Maurizio Proietti, Hirohisa Seki (Editors)

Logic-Based Program Synthesis and Transformation

Preliminary Proceedings
of the 24th International Symposium, LOPSTR 2014
Canterbury (UK), September 9–11, 2014
Preface


The aim of the LOPSTR series is to stimulate and promote international research and collaboration on logic-based program development. LOPSTR is open to contributions in all aspects of logic-based program development, all stages of the software life cycle, and issues of both programming-in-the-small and programming-in-the-large. LOPSTR traditionally solicits contributions, in any language paradigm, in the areas of synthesis, transformation, specification, analysis and verification, specialization, testing and certification, composition, program/model manipulation, optimization, transformational techniques in SE, inversion, applications and tools.

LOPSTR has a reputation for being a lively, friendly forum for presenting and discussing work in progress. Formal proceedings are produced only after the symposium so that authors can incorporate this feedback in the published papers. The formal post-conference proceedings will be published by Springer in the Lecture Notes in Computer Science series.

In response to the call for papers, 34 contributions were submitted from 21 different countries. Each submission was reviewed by at least two, and three on average, Program Committee members or external referees. The Program Committee decided to accept 20 papers for presentation at the symposium, basing this choice on their scientific quality, originality, and relevance. In addition to the 20 contributed papers, this volume includes the abstracts of the invited talks by two outstanding speakers: Roberto Giacobazzi (University of Verona, Italy), shared with PPDP, and Viktor Kuncak (EPFL, Switzerland).

We want to thank the Program Committee members, who worked diligently to produce high-quality reviews for the submitted papers, as well as all the external reviewers involved in the paper selection. We are very grateful to the LOPSTR 2014 Symposium Co-Chairs, Olaf Chitil and Andy King, and the local organizers for the great job they did in preparing the symposium. Many thanks go to Olivier Danvy, the Program Committee chair of PPDP, with whom we often interacted for coordinating the two events. Special thanks to Emanuele De Angelis and Fabrizio Smith, who helped us in maintaining the LOPSTR website and editing these proceedings. Finally, we would like to thank Andrei Voronkov
for his excellent EasyChair system that automates many of the tasks involved in chairing a conference.

Canterbury, United Kingdom
September 2014
Maurizio Proietti and Hirohisa Seki
(LOPSTR 2014 Program Co-Chairs)

Maurizio Proietti has been financially supported by the National Group of Computing Science (GNCS-INDAM).
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Invited Speakers
Obscuring Code
Unveiling and Veiling Information in Programs

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In this extended abstract we survey the most recent developments in code obfuscation and protection from a programming languages perspective. Starting from known impossibility results on universal and general purpose code obfuscation, we show that provably secure obfuscation can be achieved by constraining the attack model. This corresponds to associate attacks with suitable forms of interpretation. In this context it is always possible to systematically making code obscure, making this interpretation failing in extracting (attacking) code. The code transformation can itself be specified as the specialization of a distorted interpreter.

An extended version of this abstract is in 16th International Symposium on Principles and Practice of Declarative Programming (PPDP 2014).
Synthesizing Functions from Relations in Leon

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Abstract. We present the synthesis functionality of the Leon system (leon.epfl.ch). Leon accepts a purely functional subset of Scala extended with a choice construct. We describe automated and manual synthesis and transformation techniques in Leon, which can eliminate the choice construct and thus transform input/output relation specifications into executable functions from inputs to outputs. The techniques employed include functional synthesis procedures for decidable theories such as term algebras and Presburger arithmetic, synthesis proof rules for decomposing specifications, as well as search-based techniques, such as counterexample-guided synthesis.

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Contributed Papers
Analyzing array manipulating programs by program transformation

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Abstract. We explore a transformational approach to the problem of verifying simple array-manipulating programs. Traditionally, verification of such programs requires intricate analysis machinery to reason with universally quantified statements about symbolic array segments, such as “every data item stored in the segment $A[i]$ to $A[j]$ is equal to the corresponding item stored in the segment $B[i]$ to $B[j]$.” We define a simple abstract machine which allows for set-valued variables and we show how to translate programs with array operations to array-free code for this machine. For the purpose of program analysis, the translated program remains faithful to the semantics of array manipulation. Based on our implementation in LLVM, we evaluate the approach with respect to its ability to extract useful invariants and the cost in terms of code size.

1 Introduction

We revisit the problem of automated discovery of invariant properties in simple array-manipulating programs. The problem is to extract interesting properties of the contents of one-dimensional dynamic arrays (by dynamic we mean arrays whose bounds are fixed at array variable creation time, but not necessarily at compile time). We follow the array partitioning approach proposed by Gopan, Reps, and Sagiv [9] and improved by Halbwachs and Péron [11]. This classical approach uses two phases. In a first phase, a program analysis identifies all (potential) symbolic segments by analyzing all array accesses in the program. Each segment corresponds to an interval $I_k$ of the array’s full index domain, but its bounds are symbolic, that is, bounds are index expressions. For example, the analysis may identify three relevant segments $I_1 = [0, \ldots, i - 1]$, $I_2 = [i]$, and $I_3 = [i + 1, \ldots, n - 1]$. After this the original array $A$ is considered partitioned into segments $A_{I_k}$ corresponding to the identified segments and each segment is replaced with a summary variable $a_k$. In the second phase, the analysis aims at discovering properties $\psi(a_k)$ on each summary variable $a_k$ such that

\[ \forall \ell \in I_k (\psi(a_k) \Rightarrow \psi(A[\ell])) \] (1)
By partitioning arrays into segments, the analysis can produce stronger separate analyses for each segment rather than a single weaker combined result for the whole array. In particular, we can identify singleton segments ($A_t$ in the example) and translate array writes to these as so-called strong updates. A strong update benefits from the fact that the old content of the segment is eliminated by the update, so the new content replaces the old. For a segment that may contain multiple elements, an assignment to an array cell may leave some content unchanged, so a weak update must be used, that is, we must use a lattice-theoretic “join” of the new result and the old result associated with $\ell$.

Although very accurate, array partitioning methods have their drawbacks. Partitioning can be prohibitively expensive, with a worst-case complexity of $O(n!)$, where $n$ is the number of program variables. Moreover, partitioning must be done before the array content analysis phase that aims at inferring invariants for the form (1), which could be less precise than doing both simultaneously [5]. To mitigate this problem, the index analysis, used to infer the relevant symbolic intervals, is run twice: once during the segmentation phase and again during the array content analysis, which needs it to separate the first fixed point iteration from the rest. In the more sophisticated approach of Halbwachs and Péron [11, 16], the transfer functions are much more complex and a concept of “shift variables”, representing translation (in the geometric sense) of segments. This is not easily implemented using existing abstract interpretation libraries.

Contribution. We present a program transformation that allows scalar analysis techniques to be applied to array manipulating programs. As in previously proposed array analyses [9, 11, 16], we partition arrays into segments whose contents are treated as sets rather than sequences. To maintain the relationship among corresponding elements of different arrays, we abstract the state of all arrays within a segment to a set of vectors, one element per array. Thus we transform an array manipulating program into one that manipulates scalars and sets of vectors. A major challenge in this is to encode the disjunctive information carried by each array segment. We propose a technique that splits basic blocks. It has been implemented using the LLVM framework.

Importantly, a program transformation approach allows the separation of concerns: existing analyses based on any scalar abstract domains can be used directly to infer array content properties, even interprocedurally. While other approaches lift a scalar abstract domain to arrays by lifting each transfer function, our approach uses existing transfer functions unchanged, only requiring the addition of two simple transfer functions easily defined in terms of operations that already exist for most domains. The approach is also parametric in the granularity of array index sets, ranging from array smashing [2] to more precise (and expensive) instances. When we go beyond array smashing, the transformational approach inherits the exponential search cost present in the Halbwachs/Péron approach, as for some programs $P$, the transformed programs $P'$ are exponentially larger than $P$. However, for simple array-based sort/search programs [9, 11], a transformational approach is perfectly affordable, in particular as we can capitalize on code optimization support offered by the LLVM infrastructure.
Analyzing array manipulating programs by program transformation

Instructions
\[ I \rightarrow v_1 = \text{constant} \mid v_1 = \circ v_2 \mid v_1 = v_2 \circ v_3 \mid A \]

Array assignments
\[ A \rightarrow v_1 = \text{arr}[v_2] \mid \text{arr}[v_1] = v_2 \]

Jumps
\[ J \rightarrow \text{If} (v_1 \circ v_2) \écial \text{label}_1 \text{ label}_2 \mid \text{Jmp} \text{ label} \mid \text{error} \mid \text{end} \]

Blocks
\[ B \rightarrow \text{label} : \text{I}^* \text{J} \]

Programs
\[ P \rightarrow B^+ \]

---

Fig. 1. A small control-flow language with array expressions

---

2 Source and target language

Our implementation uses LLVM IR as source and target language. However, as the intricacies of static single-assignment (SSA) form obscure, rather than clarify, the transformation, we base our presentation on a small traditional control flow language, whose syntax is given in Fig. 1. We shall usually shorten “basic block” to “block” and refer to a block’s label as its identifier.

**Source.** Each block is a (possibly empty) sequence of instructions, followed by a (conditional) jump. Arithmetic unary and binary operators are denoted by \( \circ \) and \( \diamond \) respectively, and logical operators by \( \triangledown \). We assume that there is a fixed set of arrays \( \{A_1, \ldots, A_k\} \), which have global scope (and do not overlap in memory). The semantics is conventional and not discussed here. Fig. 2(a) shows an example program in diagrammatic form.

**Target.** The abstract machine we consider operates on variables over two kinds of domains: standard scalar types, and sets of vectors of length \( k \), where \( k \) is the number of arrays in the source. The scalar variables represent scalars of the source program, including index variables, as well as singleton array segments; sets of vectors represent non-singleton segments of all extant arrays. Let \( \mathbf{V} \) be the
set of scalar variables and $S$ be the set of vector set variables. The runtime state of the machine is given by a pair $\langle \sigma, \rho \rangle$ consisting of a variable store $\sigma : V \to \mathbb{Z}$, and a set store $\rho : S \to \mathcal{P}(\mathbb{Z}^k)$.

A control flow language for set machine programs is given in Fig. 3. Arithmetical and logical operations affect only the variable store $\sigma$; the semantic rules for these operations are standard. The set machine also has set operations union ($\cup$), subset (nondet-subset) and element of (nondet-elt). Fig. 4 gives their semantic rules, distinguishing scalar variables $v$ and (vector) set variables $S$.

The union update $S_1 = S_2 \cup S_3$ maps $S_1$ to the union of values of $S_2$ and $S_3$. The subset and element operations are non-deterministic: executing $S_1 = \text{nondet-subset}(S_2)$ assigns to $S_1$ some non-empty subset of elements from $S_2$, but makes no guarantee as to which elements are selected. Similarly, the element operation $(v_1, \ldots, v_k) = \text{nondet-elt}(S_1)$ nondeterministically selects some element of vector set $S_1$ to load into $v_1, \ldots, v_k$.

**Translation.** Fig. 2(a)’s program scans an array for the first occurrence of value $x$, assumed to occur in $A$. The constraint $A[i] = x \wedge \forall k \in [0, i)$ ($A[k] \neq x$) is the desired invariant at tail. A corresponding array-free program is given in Fig. 2(b). The example illustrates some key features. Each contiguous array segment is represented by a set variable $A_i$. Each original block is duplicated for each feasible ordering of interesting variables. In the initial ordering ($0 = i$) the only interesting segment is $A[0]$, represented as the singleton $a_0$; the read $v = A[i]$ is replaced by an assignment $v = a_0$. At guard$^{0<i}$, $A[i]$ is represented by $a_1$ so the read is replaced by $v = a_1$. When $i$ is updated at body$^{0=i}$, the previous singleton $a_0$ becomes part of an “aggregate” segment $A[0, i-1]$. We then transform singleton $a_0$ to set $A_0$ and introduce a new singleton $a_1$ (representing $A[i]$ in the updated ordering). Similarly, when we update $i$ in body$^{0<i}$, segments $A[0, i-1]$ and $A[i]$ are merged (yielding $A_0 = A_0 \cup \{a_1\}$), and a new singleton $a_1$ is introduced. Consider the resulting concrete set-machine states. At tail$^{0<i}$, we have $a_0 = x$, corresponding to $A[0] = x$ in the original program. At tail$^{0=i}$, we find $x \notin A_0$ and $a_1 = x$. These correspond, respectively, to array invariants $\forall \ell \in [0, i-1] : A[\ell] \neq x$ and $A[i] = x$ in the original program.
3 From scalar to set machine transfer functions

We now show how to lift a scalar domain for use by set machines. Essentially, we use a scalar variable to approximate each component of each set; approximation of set-machine states can then be obtained by grouping states by values of the (original) scalar variables. Essentially, we approximate a set-machine state \(\langle \sigma, \rho \rangle\) with set variables \(\{S_1, \ldots, S_m\}\) by a set of scalar states \(\{S_1^\circ, \ldots, S_m^\circ\}\) representing the possible results of selecting some element from each set:

\[
\alpha^\circ(\langle \sigma, \rho \rangle) = \{\sigma \cup \{S_1^\circ \mapsto y_1, \ldots, S_m^\circ \mapsto y_m\} \mid y_1 \in \rho(S_1), \ldots, y_m \in \rho(S_m)\}
\]

Transfer functions for set operations then operate over the universe of possible states, rather than apply element-wise to each state.

Example 1. Consider a program with one scalar variable \(x\), and one set variable \(S\), with initial state \(\langle\{x \mapsto 0\}, \{S \mapsto \{1, 2\}\}\rangle\). If we introduce a scalar variable \(y\) that selects a value from \(S\) via \(a\) we have two possible states:

\[
\langle\{x \mapsto 0, y \mapsto 1\}, \{S \mapsto \{1, 2\}\}\rangle, \langle\{x \mapsto 0, y \mapsto 2\}, \{S \mapsto \{1, 2\}\}\rangle
\]

If we represent \(S\) by scalar variable \(S^\circ\), we have initial states \(\langle\{x \mapsto 0, S^\circ \mapsto \{1\}\}\rangle\) and \(\langle\{x \mapsto 0, S^\circ \mapsto \{2\}\}\rangle\). When we wish to select a value for \(y\), it is chosen nondeterministically from the possible values of \(S^\circ\), resulting in the states:

\[
\langle\{x \mapsto 0, y \mapsto 1, S^\circ \mapsto \{1\}\}\rangle, \langle\{x \mapsto 0, y \mapsto 1, S^\circ \mapsto \{2\}\}\rangle
\]

\[
\langle\{x \mapsto 0, y \mapsto 2, S^\circ \mapsto \{1\}\}\rangle, \langle\{x \mapsto 0, y \mapsto 2, S^\circ \mapsto \{2\}\}\rangle
\]

If we group states with equal values of \(x\) and \(y\), we can see that these correspond to the final states of the original set-machine fragment.

Note that this is an (over-)approximation. We can only infer the set of values that may be elements of \(S\)—this representation cannot distinguish sets of elements which may occur together, nor the cardinality of \(S\). For example, assume we have possible set-machine states \(\emptyset, \{S \mapsto \{1\}\}\) and \(\emptyset, \{S \mapsto \{2\}\}\). The scalar approximation is \(\langle\{S^\circ \mapsto \{1\}\}, \{S^\circ \mapsto \{2\}\}\rangle\) which covers the feasible set-machine states, but also includes \(\emptyset, \{S \mapsto \{1, 2\}\}\), which is not feasible. More generally, if the set \(\varphi\) of concrete states allows sets \(S \mapsto X_1, \ldots, S \mapsto X_k\), we have:

\[
\forall X_1 \ldots X_k \ (X_1 \subseteq X \cup \ldots \cup X_k \Rightarrow (S \mapsto X) \in \gamma \circ \alpha(\varphi))
\]

Consider a (not necessarily numeric) abstract domain \(A\), with meet (\(\cap\)), join (\(\cup\)) and rename operations, as well as a transfer function \(\mathcal{F} : I \rightarrow A \rightarrow A\) for the scalar fragment of the language. The rename operation constructs a new state where each variable \(x_i\) is replaced with \(y_i\) (then removes the existing bindings of \(x_i\)). Formally, the concrete semantics of rename is given by

\[
\text{rename}(\sigma, [x_1, \ldots, x_k], [y_1, \ldots, y_k]) = \sigma\left[\begin{array}{c}
y_1 \mapsto \sigma(x_1), \ldots, y_k \mapsto \sigma(x_k), \\
x_1 \mapsto *, \ldots, x_k \mapsto *
\end{array}\right]
\]

For each set variable \(S\), we introduce \(k\) scalar variables \([s^1, \ldots, s^k]\) denoting the possible values of each vector in \(S\). We then extend \(\mathcal{F}\) to set operations as shown in Fig. 5.
we will denote by $I$ a total ordering of a set of index variables $B$.

For convenience in expressing these orderings, we will introduce for each index variable $i<j$ a removal followed by an introduction.

$k$ and $j$ leaves scope, computing the resulting ordering consists simply of eliminating $k$ and $j$ from $\pi$, and discarding any now-empty sets. Assignment of an index variable is handled as a removal followed by an introduction.\footnote{3}

\begin{align*}
\mathcal{F}[S_1 = \text{nondet-subset}(S_2)] \varphi &= \varphi \cap \text{rename}(\varphi, [s_1^2, \ldots, s_2^2], [s_1^1, \ldots, s_2^1]) \\
\mathcal{F}[S_1 = S_2 \cup S_3] \varphi &= \left( \mathcal{F}[S_1 = \text{nondet-subset}(S_2)] \varphi \cup \mathcal{F}[S_1 = \text{nondet-subset}(S_3)] \varphi \right) \\
\mathcal{F}[(v_1, \ldots, v_k) = \text{nondet-elt}(S_1)] \varphi &= \varphi \cap \text{rename}(\varphi, [s_1^1, \ldots, s_k^1], [v_1^1, \ldots, v_k^1])
\end{align*}

Fig. 5. Extending the transfer function for scalar analysis to set operations.

4 Orderings

The transformation relies on maintaining a single ordering of index variables at each transformed block. We now discuss such total orderings.

Our goal is to partition the array index space $(-\infty, \infty)$ into contiguous regions bounded by index variables. For index variables $i$ and $j$, we need to be able to distinguish between the cases where $i < j$, $i = j$, and $i > j$. However, this is not enough; if we assign $A[i] = x$, but only know that $i < j$, we cannot distinguish between the cases $i = j - 1$ (every element between $i$ and $j$ is $x$) and $i < j - 1$ (there are additional elements with some other property). So, for index variables $i$ and $j$, we distinguish these five cases:

\[
\begin{array}{cccc}
i + 1 < j & i + 1 = j & i = j & i = j + 1 \\
i > j + 1
\end{array}
\]

For convenience in expressing these orderings, we will introduce for each index variable $i$ a new term $i^+$ denoting the value $i + 1$, and for a set of index variables $I$ we will denote by $I^+$ the augmented set $I \cup \{v^+ \mid v \in I\}$. We can then define a total ordering of a set of index variables $I$ to be a sequence of sets $[B_1, \ldots, B_k]$, $B_s \subseteq I^+$, such that the $B_s$’s cover $I^+$, are pairwise disjoint, and satisfy $i \in B_s \Leftrightarrow i^+ \in B_{s+1}$.

The meaning of the ordered list $\pi = [B_1, B_2, \ldots, B_k]$ is parameterised by the value of program variables involved, that is, it depends on a store $\sigma$. The meaning is: $[\pi](\sigma) \equiv $

\[
\bigwedge_{s, t \in [1..k]} (\forall e, e' \in B_s \ (\sigma(e) = \sigma(e')) \land \forall e \in B_s \ (\forall e' \in B_t \ (s < t \rightarrow \sigma(e) < \sigma(e'))))
\]

An ordering $\pi$ (plus virtual bounds $(-\infty, \infty)$) partitions the space of possible array indices into contiguous regions, given by $[\sigma(e), \sigma(e')]$ for $e \in B_s, e' \in B_{s+1}$.

For any index variable $i$, a segment containing $i^+$ in the right bound is necessarily a singleton segment; all other segments are considered aggregate.

When a new index variable $k$ enters scope, several possible orderings may result. Fig. 6(c) gives a procedure for enumerating them. When an index variable $k$ leaves scope, computing the resulting ordering consists simply of eliminating $k$ and $k^+$ from $\pi$, and discarding any now-empty sets. Assignment of an index variable is handled as a removal followed by an introduction.\footnote{3}
We can discard any ordering that arranges constants in infeasible ways, such as $4 < 3$. If we have performed some scalar analysis on the original program, we need only generate orderings which are consistent with the analysis results.

5 The transformation

We now detail the transformation from an array manipulating program to a set-machine program, with respect to a fixed set of interesting segment bounds. Section 6 covers the selection of these bounds. Intuitively, the goal of the transformation is to partition the array into a collection of contiguous segments, such that each array operation uniquely corresponds to a singleton segment. Each singleton segment is represented by a tuple of scalars; each non-singleton segment is approximated by a set variable. There are two major obstacles to this. First, a program point does not typically admit a unique ordering of a given set of segment bounds; second, as variables are mutated in the program, the correspondence between concrete indices and symbolic bounds changes.

The transformation resolves this by replicating basic blocks to ensure that, at any program point, a unique partitioning of the array into segments is identifiable. Any time a segment-defining variable is modified, introduced or eliminated, we emit statements to distinguish the possible resulting partitions, and duplicate the remainder of the basic block for each case. For each partition, we also emit set operations to restore the correspondence between set variables and array segments, using \texttt{nondet-el} and \texttt{nondet-subset} when a segment is subdivided, and $\cup$, when a boundary is removed, causing segments to be merged. This way every array read/write in the resulting program can be uniquely identified with a singleton segment. As singleton sets are represented by tuples of scalars, we can finally eliminate array operations, replacing them with scalar assignments.

In the following, we assume the existence of functions \texttt{next-block}, which allocates a fresh block identifier, and \texttt{push-block}, which takes an identifier, a sequence of statements and a branch, and adds the resulting block to the program. We also assume that there is a mutable global table $T$ mapping block identifier and index variable ordering pairs $\langle id, \pi \rangle$ to $ids$, used to store previously computed partial transformations, and an immutable set $I$ of segment bound variables and constants. The function \texttt{get-block} takes a block identifier, and returns the body of the corresponding block. The function \texttt{vars} returns the set of variables appearing lexically in the given expression. The function \texttt{find-avar} gives the variable name to which a given array and index will be translated, given an ordering.

Fig. 6 gives the transformation. Procedure \texttt{transform} takes a block and transforms it, assuming a given total ordering $\pi$ of the index variables. It is called once with the initial block of each function and an ordering containing only the constants in the index set. As there are finitely many $\langle id, \pi \rangle$ combinations, and each pair is constructed at most once, this process terminates.

The core of the transformation is done by a call to \texttt{transform-body}(B, $\pi$, id, ss). Here $B$ is the portion of the current block to be transformed and $\pi$ the current ordering. id and ss hold the identifier and body of the partially-transformed
% Check if the block has already been transformed
% under π. If not, transform it.
transform(id, π)
  if (⟨id, π⟩ ∈ T)
    return T[⟨id, π⟩]
  idı := next_block()
  T := T[⟨id, π⟩ → idı]  
  (stmts, br) := get_block(idı)
  transform_body((stmts, br), π, id, [])
  return idı

% Evaluate a branch.
transform_body([], Jmp b, π, id, ss)
  idb := transform(b, π)
  push_block(id, ss, Jmp idb)

transform_body([key where l then t else f], π, id, ss)
  if vars(l) ⊆ I
    dest := if eval(l, π) then t else f
    iddest := transform(dest, π)
    push_block(id, ss, Jmp iddest)
  else
    idı := transform(t, π)
    idf := transform(f, π)
    push_block(id, ss, If l then idı else idf)
  % (Potentially) update an index.
  transform_body([key where x = expr[stmts], br], π, id, ss)
  if x ∈ I
    split_transform(x, (stmts, br), π, id, ss++x = expr)
  else
    transform_body((stmts, br), π, id, ss++x = expr)

% Transform an array read...
transform_body([key where A[i] = x[stmts], br], π, id, ss)
  Ai := find_avar(π, A, i)
  transform_body((stmts, br), π, id, ss++A = Ai)
% or an array write.
transform_body([key where A[i] = x[stmts], br], π, id, ss)
  Ai := find_avar(π, A, i)
  transform_body((stmts, br), π, id, ss++A[i] = x)

(a) The top-level transformation process

split_transform(x, (stmts, br), π, id, ss)
  II′ := feasible_orders(π, x)
  split_rec(x, II′, (stmts, br), π, id, ss)

split_rec(x, [π′], (stmts, br), π, id, ss)
  asts := remap_avars(π, π′)
  transform_body((stmts, br), π′, id, ss++asts)

split_rec(x, [π′|II′], (stmts, br), π, id, ss)
  idast := next_block()
  idII′ := next_block()
  cond := ord_cond(x, π′)
  push_block(id, ss, If cond then idast else idII′)
  asts := remap_avars(π, π′)
  transform_body((stmts, br), π′, idast+asts)
  split_rec(x, II′, (stmts, br), π, idII′, [])

(b) Fan-out of a block when an index variable is changed

feasible_orders(k, π) := insert(k, π, [])

insert(k, [], pre) := return {pre++k : {k}}

insert(k, [S] | S, pre)
  low := insert*(k, S, [S], pre++{k})
  high := insert(k, S, pre++{S})
  if ∃ x . x ∈ S
    return low ∪ high
  else
    return low ∪ high
      insert*(k, S, pre++{S ∪ {k}})

insert*(k, [], pre) := return {pre++{k}}

insert*(k, [S] | S, pre)
  if ∃ x . x ∈ S
    return {pre++{S ∪ {k}} : {S}}
  else
    return {pre++{S ∪ {k}} : {S} ∪
             {pre++{k}} : {S : S}}

(c) Enumerating the possible total orderings upon introducing a new index variable k

Fig. 6. Pseudo-code for stages of the transformation process.
Analyzing array manipulating programs by program transformation

\[ \pi = \{0\} < \{1\} < \{i\} < \{i^+\} < \{n\} \]

\[
\begin{align*}
\text{transform}_\text{body} : \pi & \quad x := a_2 \\
B[i] := x & \quad i := i + 1 \\
\text{transform}_\text{body} : \pi & \quad b_2 := x \\
i := i + 1 & \\
\end{align*}
\]

(a) Original (b) After Step 1 (c) After Step 2

Fig. 7. Transformation of array reads and writes under ordering \(\pi\). As the segment \([i, i^+]\) is a singleton, the array elements are represented as scalars.

accesses and assignments to the corresponding scalar or set variable, as determined by find_aivar. Conditional branches whose conditions are determined by the current ordering are replaced by direct branches to the then or else part, as appropriate. Once no instructions remain to be transformed, the block id is emitted with body ss, together with the appropriate branch instruction.

Whenever an index variable is modified, the rest of the current block must be split, and the set variables must be updated accordingly. The rest of the block is then transformed under each possible new ordering \(\pi'\). This is the job of split_transform shown in Fig. 6(b), while the job of feasible_orders in Fig. 6(c) is to determine the set of possible orders. The function ord_cond\((x,\pi')\) generates logical expressions to determine whether the ordering \(\pi'\) holds, given that \(\pi\) previously held. ord_cond checks the position of both \(x\) and \(x^+\). If \(x\) is part of a larger equivalence class in \(\pi\), ord_cond generates the corresponding equality; otherwise, it checks that \(x\) is greater than its left neighbour; similarly, it checks that \(x^+\) is in its class or less than its right neighbour. Fig. 6(b) shows the process of splitting a block upon introducing an index variable \(x\).

5.1 Reading and writing

Transformation of array reads and writes is simple, if the array index is in the set \(I\) of index variables. Fig. 7(a–c) shows the step-by-step transformation of a block, under the specified ordering. After Step 1, reference \(A[i]\) has been transformed to scalar \(a_2\), since \(\{i\}\) is a singleton. Similarly, Step 2 transforms \(B[i]\) to \(b_2\).

If the index of the read/write operation has been omitted, we must instead emit code to ensure the operation is dispatched to the correct set variable. The dispatch procedure is similar in nature to split_transform, as given in Fig. 6(c); essentially, we emit a series of branches to determine which (if any) of the current segments contains the read/write index. Once this has been determined, we apply the array operation to the appropriate segment. If the selected segment is a singleton, this is done exactly as in transform_body. For writes to an aggregate segment, we must first read some vector from the segment, substitute the element to be written, then merge the updated vector back into the segment4.

4 Detailed pseudo-code for this is in Appendix A.
\( \pi = [\{0\} < \{1\} < \{i\} < \{i^+\} < \{n\}] \)
\( \pi'_0 = [\{0\} < \{1\} < \{i\} < \{i^+\} < \{n\}] \)
\( \pi'_1 = [\{0\} < \{1\} < \{i\} < \{i^+, n\}] \)

\[\begin{align*}
\text{transform} &\ : \pi \\
i &:= i + 1 \\
\text{Jmp guard} &
\end{align*}\]

(a) Original

\[\begin{align*}
i &:= i + 1 \\
\text{split} &\ : \pi \mapsto [\pi'_0, \pi'_1] \\
\text{Jmp guard} &
\end{align*}\]

(b) After Step 1

\[\begin{align*}
\text{remap avars} &\ : \pi \mapsto \pi'_0 \\
\text{transform} &\ : \pi'_0 \\
\text{Jmp guard} &
\end{align*}\]

\[\begin{align*}
s_{\pi'_0} &
\end{align*}\]

(c) After Step 2

\[\begin{align*}
\text{remap avars} &\ : \pi \mapsto \pi'_1 \\
\text{transform} &\ : \pi'_1 \\
\text{Jmp guard} &
\end{align*}\]

(d) After Step 3

Fig. 8. Example of updating an index assignment. We assume an existing scalar analysis which has determined that, after \( i = i + 1 \), we have \( 1 < i < n \).

5.2 Index manipulation

The updating of index variables is the most involved part of the transformation, as we must emit code not only to determine the updated ordering \( \pi' \), but also to ensure the array segment variables are matched to the corresponding bounds. Fig. 8 illustrates this process, implemented by the procedure \text{remap avars}, as it splits a block into three: one to test an index expression to determine what ordering applies, and one for each ordering. In the original code, ordering \( \pi \) applies, but following the assignment, either ordering \( \pi'_0 \) or \( \pi'_1 \) may apply. The test inserted by Step 2 distinguishes these cases, leaving only one ordering applicable to each of the \( s_{\pi'_0} \) and \text{split} \_1 \) blocks.

If we normalize index assignments such that for \( k := E, k \notin E \), we can separate the updating of segment variables into two stages; first, computing intermediate segment variables \( A'_i \) after eliminating \( k \) from \( \pi \), and then computing the new segment variables after introducing the updated value of \( k \). Pseudo-code for these steps are given in Fig. 9(a) and 10(a). In practice, we can often eliminate many of these intermediate assignments, as segments not adjacent to the initial or updated values of \( k \) remain unchanged.
Analyzing array manipulating programs by program transformation

remap_avars(k, π, π')
   eliminate(k, π) :: introduce(k, π')

eliminate(k, π)
   eliminate(k, π, 0, 0, 0)

eliminate(k, [], i', E)
   if(i' = 0) return []
   else return [emit_merge(A'_{i'-1}, E)]

eliminate(k, [c] | S, i, i', E)
   where c ∈ {k, k^+}
   return eliminate(k, S, i + 1, i', E ∪ {A_i})

eliminate(k, [S_j] | S, i, i', E)
   suff := eliminate(k, S, i + 1, i' + 1, {A_i})
   if(i' = 0)
      % Ignore leading segment.
      return suff
   else
      return emit_merge(A'_{i'-1}, E) :: suff

emit_merge(x, E)
   return [x = \bigcup E]

(a)

\[
\begin{align*}
\pi &= [(k) < \{k^+, n\} < \{n^+\}] \\
p_r &= [(n) < \{n^+\}] \\
A_0 &:= A_1
\end{align*}
\]

(b)

\[
\begin{align*}
\pi &= [(k, n) < \{k^+, n^+\}] \\
p_r &= [(n) < \{n^+\}] \\
A_0 &:= A_0
\end{align*}
\]

\[
\begin{align*}
\pi &= [(i, k) < \{k^+, n\} < \{n^+\}] \\
p_r &= [(i) < \{n^+\}] \\
A_0 &:= A_1 \\
A_1 &:= \{A_1\} \cup A_2 \\
A_2 &:= A_3
\end{align*}
\]

Fig. 9. (a) Algorithm for generating instructions to keep segment variables updated;
(b) resulting assignments when k is eliminated from various orderings, also showing
the remaining order \(p_r\) and scalar or set variables corresponding to each segment.

When we eliminate an index variable \(k\) from \(\pi\), we merge segments that were
bounded only by \(k\) or \(k^+\). If \(k\) or \(k^+\) appears alone at the very beginning or
end of \(\pi\), the segments are discarded entirely. If either appears alone between
other variables in \(\pi\), the segments on either side are merged to form a single
segment. However, if \(k\) and \(k^+\) are both equal to some other variables, the
original segments are simply copied to the corresponding temporary variables.
This is illustrated in Fig. 9(b).

The pseudo-code in Fig. 9 and 10 ignores the distinction between singleton
and aggregate segments; the transformed operations differ slightly in the two
cases. If we introduce a singleton segment into an aggregate segment, we select
a single vector from the set \([(a', b', c') = \text{nondet-elt}(A_i)]; if an aggregate segment
is introduced, we emit a subset operation \((A' = \text{nondet-subset}(A))\).

The procedure for injecting \(k\) into \(\pi\) behaves similarly. If \(k\) is introduced
at either end of \(\pi\), we introduce new segments with indeterminate values. If \(k\)
introduce\((k, \pi)\)
introduce\((k, \pi, 0, 0)\)
introduce\((k, [], i, i')\)
\text{return} []
introduce\((k, [c] \mid S], i, i')\)
\text{where} \(c \in \{k, k^+\}\)
\text{suff} := introduce\((k, S, i + 1, i')\)
\text{if} \((i' = 0)\)
\text{return} \([A_i = *] : : \text{suff}\)
\text{else}
\text{return} \([A_i = \text{nondet-subset}(A'_{i'])] : : \text{suff}\)
introduce\((k, [S_j \mid S], i, i')\)
\text{suff} := introduce\((k, S, i + 1, i' + 1)\)
\text{if} \((i' = 0)\)
\text{\% Ignore leading segment.}
\text{return} \text{suff}
\text{else}
\text{return} \([A_i = A'_{i'}] : : \text{suff}\)

\(\pi_p = [\{n\} < \{n^+\}]\)
\(\pi = [\{k\} < \{k^+, n\} < \{n^+\}]\)
\(\begin{align*}
a_0 & := * \\
a_1 & := a'_0 \\
\end{align*}\)
\(\pi_p = [\{n\} < \{n^+\}]\)
\(\pi = [\{k, n\} < \{k^+, n^+\}]\)
\(\begin{align*}
a_0 & := a'_0 \\
\end{align*}\)
\(\pi_p = [\{i\} < \{i^+\} < \{n\} < \{n^+\}]\)
\(\pi = [\{i\} < \{i^+, k\} < \{k^+\} < \{n\} < \{n^+\}]\)
\(\begin{align*}
a_0 & := a'_0 \\
(a_1) & = \text{nondet-elt}(A'_{i'}) \\
A_2 & = \text{nondet-subset}(A'_{i'}) \\
A_3 & := a'_2 \\
\end{align*}\)

Fig. 10. (a) Generating instructions for (re-)introducing a variable \(k\) into a given ordering, and (b) the resulting assignments when \(k\) is introduced into various orderings. Note the difference between introducing singleton and aggregate segments.

is introduced somewhere within an existing segment, we introduce new child segments—each of which is a subset of the original segment.

5.3 Control flow

When transforming control flow, there are three cases we must consider:

1. Unconditional jumps
2. Conditional jumps involving some non-index variables
3. Conditional jumps involving \textit{only} index variables

In cases (1) and (2), the transformation process operates as normal; we recursively transform the jump targets, and construct the corresponding jump with the transformed identifiers. However, when we have a conditional jump \textbf{If} \(i \bowtie j\) \textbf{then} \(t\) \textbf{else} \(f\) where \(i\) and \(j\) are both index terms, the relationship between \(i\) and \(j\) is statically determined by the current ordering \(\pi\). As a result, we can simply evaluate the condition \(i \bowtie j\) under the ordering \(\pi\), and use an unconditional branch to the corresponding block. This is illustrated in Fig. 11.
Analyzing array manipulating programs by program transformation

\[
\pi = [\{0\} < \{1\} < \{i\} < \{i^+\} < n]
\]

\[
\text{let } body_\pi = \text{transform}(body, \pi) \text{ in }
\]

\[
\text{transform}_\pi : \pi
\]

\[
\begin{align*}
\text{If } i < n & \quad \text{then } body \\
\text{else } tail
\end{align*}
\]

\[
\text{Jmp } body_\pi
\]

Fig. 11. Transforming a jump, conditional on index variables only, under ordering \(\pi\)

6 Selecting segment bounds

In the previous sections, we have assumed a pre-selected set of interesting segment bounds. The selection of segment boundaries involves a trade-off: we can improve the precision of the analysis by introducing additional segment bounds, but the transformed program grows exponentially as the number of segment bounds increases. As do [11], we can run a data-flow analysis to find the set of variables that may (possibly indirectly) be used as, or to compute, array indices. Formally, we collect the set \(I\) of variables and constants \(i\) occurring in these contexts:

- \(A[i]\) where \(A\) is an array (2)
- \(i' = \text{op}(i)\) where \(i' \in I\) (3)
- \(i = \text{op}(i')\) where \(i' \in I\) and \(i'\) is not a constant (4)

Any variable which does not satisfy these conditions can safely be discarded as a possible segment bound. For the experiments in Section 7 we used all elements of \(I\) as segment bounds (so \(I = I\)), which yields an analysis corresponding roughly to the approaches of [9, 11]. We could, however, discard some subset of \(I\) to yield a smaller, but less precise, approximation of the original program. The cases (3) and (4) are needed because of possible aliasing; this is particularly critical in an SSA-based language, as SSA essentially replaces mutation with aliasing. It is worth noting that these dependencies extend to aliases introduced prior to the relevant array operation, as in the snippet “\(i := x; \ldots A[i] := k; \ldots y := x + 1;\)”

7 Experimental evaluation

We have implemented our method using the LLVM framework, with the transformation carried out over two distinct passes. In the first pass, transformation is done as described above, but without great regard for the number of basic blocks and instructions in the transformed program. At the same time, we also use a scalar analysis of the original program to detect any block whose total ordering is infeasible. In the second pass, we prune these unreachable blocks away. As can be gleaned from Table 1, these measures reduce the complexity of the transformed program significantly.
In order to infer array properties from invariants discovered by analysing transformed versions of programs, we require users to specify, at transformation time, the range of array segments that they are interested in reasoning about. In our implementation, this is described by a strict inequality between indices that must hold for segments in the range. For example, specifying $0 < n$ indicates that we want to discover invariants of the form $\forall \ell (0 \leq \ell < n \Rightarrow \psi(A[\ell]))$, where $A$ is some array and $\psi$ is some property. At the end of the transformation we use a newly created block to join all copies of the original exit block whose total ordering is consistent with the given range. The various scalar representations for each array segment, as well as other variables in scope in each copy, are merged together in phi nodes inside this final block. Properties discovered about the segment phi nodes then translate directly to properties about the corresponding array segments in the original program.

We have tested our method by running first the polka polyhedra domain [12] on the output of our transformation when applied to the programs given in Fig. 12. The interesting invariants that we infer are as follows (each property holds at the end of the corresponding function):

<table>
<thead>
<tr>
<th>Function</th>
<th>Invariant</th>
</tr>
</thead>
<tbody>
<tr>
<td>array_copy (int* A, int* B, int n)</td>
<td>$\forall \ell (0 \leq \ell &lt; n \Rightarrow A[\ell] = B[\ell])$</td>
</tr>
<tr>
<td>array_init (int* A, int n)</td>
<td>$\forall \ell (0 \leq \ell &lt; n \Rightarrow A[\ell] = 5)$</td>
</tr>
<tr>
<td>array_max (int* A, int n)</td>
<td>$\forall \ell (0 \leq \ell &lt; n \Rightarrow A[\ell] \leq \text{max})$</td>
</tr>
<tr>
<td>search (int* A, int key)</td>
<td>$\forall \ell (0 \leq \ell &lt; i \Rightarrow A[\ell] \neq \text{key})$</td>
</tr>
<tr>
<td>first_not_null (int* A, int n)</td>
<td>$\forall \ell (0 \leq \ell &lt; s \Rightarrow A[\ell] = 0)$</td>
</tr>
<tr>
<td>sentinel (int* A, int n, int sent)</td>
<td>$\forall \ell (0 \leq \ell &lt; i \Rightarrow A[\ell] \neq \text{sent})$</td>
</tr>
</tbody>
</table>

Fig. 12 shows test programs from related papers [9, 11]. Table 1 lists sizes of the original, transformed, and post-processed transformed versions (columns Original, Transformed, and Post-processed respectively), as well as the time to perform the transformation (column transf). Column polka shows the analysis time in seconds for running the polka polyhedra domain. uva is explained below.
Analyzing array manipulating programs by program transformation

<table>
<thead>
<tr>
<th>Program</th>
<th>Original</th>
<th>Transformed</th>
<th>Post-processed</th>
<th>Running time (s)</th>
</tr>
</thead>
<tbody>
<tr>
<td></td>
<td>blocks</td>
<td>insts.</td>
<td>blocks</td>
<td>insts.</td>
</tr>
<tr>
<td>array_copy</td>
<td>5</td>
<td>12</td>
<td>274</td>
<td>898</td>
</tr>
<tr>
<td>array_init</td>
<td>5</td>
<td>11</td>
<td>274</td>
<td>644</td>
</tr>
<tr>
<td>array_max</td>
<td>7</td>
<td>19</td>
<td>220</td>
<td>562</td>
</tr>
<tr>
<td>search</td>
<td>5</td>
<td>10</td>
<td>90</td>
<td>167</td>
</tr>
<tr>
<td>first_not_null</td>
<td>8</td>
<td>17</td>
<td>1057</td>
<td>2217</td>
</tr>
<tr>
<td>sentinel</td>
<td>5</td>
<td>13</td>
<td>1001</td>
<td>1936</td>
</tr>
</tbody>
</table>

Table 1. Sizes of transformed test programs and analysis time for polka and uva

Enhancing an existing analyzer. As a separate experiment we use IKOS [3], an abstract interpretation-based static analyzer developed at NASA. IKOS has been used successfully to prove absence of buffer overflows in Unmanned Aircraft Systems flight control software written in C. The latest (unreleased) version of IKOS provides an uninitialized variable analysis that aims at proving that no variable can be used without being previously defined, otherwise the execution of the program might result in undefined behaviour.

Currently IKOS is not sufficiently precise for array analysis. (Fig. 13), IKOS cannot deduce that \( A[5] \) is definitely uninitialized at line 4. However, using the transformational approach, IKOS proves that \( A[5] \) is definitely uninitialized. The problem is far from trivial; as Regehr [17] notes, gcc and clang (with \(-Wuninitialized\)) do not even raise warnings for this example, but stay completely silent.

We ran IKOS on the transformed version of array_init_unsafe. IKOS successfully reported a definite error at line 4 in 0.22 seconds. Conversely, transformation enabled IKOS to show that no array element was left undefined in the case of array_init. Finally we ran IKOS on the rest of the programs in Fig. 12. For the purpose of the uninitialized variable analysis we added loops to force each array to be treated as initialized, when appropriate. For the transformed version of array_copy, IKOS proved that \( A \) is definitely initialized after the execution of the loop. For the rest of the programs IKOS proved that the initialized array \( A \) is still initialized after the loops. Column uva in Table 1 shows the analysis time in seconds of the uninitialized variable analysis implemented in IKOS.

Note that the polka analysis does not eliminate out-of-scope variables. Our program transformation introduces many variables, and since polka incurs a super-linear per-variable cost, the overall time penalty is considerable. We expect to be able to greatly reduce the cost by utilising a projection operation and improving the fixed-point finding algorithm.
8 Related work

Amongst work on automated reasoning about array-manipulating code, we can distinguish work on analysis from work that focuses on verification. Our paper is concerned with the analysis problem, that is, how to use static analysis for automated generation of (inductive) code invariants. As mentioned in Section 1, we follow the tradition of applying abstract interpretation [4] to the array content analysis problem [5, 9, 11, 16]. Alternative array analysis methods include Gulwani, McCloskey and Tiwari’s lifting technique [10] (requiring the user to specify templates that describe when quantifiers should be introduced), Kovács and Voronkov’s theorem-prover based method [13], Dillig, Dillig and Aiken’s fluid updates [7] (supporting points-to and value analysis but excluding relational analyses), and incomplete approaches based on dynamic analysis [8, 15].

Unlike previous work, we apply abstract interpretation to a transformed program in which array reads and writes have been translated away; any standard analysis, relational or not, can be applied to the resulting program, with negligible additional implementation cost.

There is a sizeable body of work that considers the verification problem for array-processing programs. Here the aim is to establish that given assertions hold at given program points. While abstract interpretation may serve this purpose (given a well-chosen abstract domain), more direct approaches are goal-directed, using assertions actively, to drive reasoning, rather than passively, as checkpoints. Many alternative techniques have been suggested for the verification of (sometimes restricted) array programs, including lazy abstraction [1], template-based methods [14], and, more closely related to the present paper, techniques that employ translation, for example to Horn clauses [6].

9 Conclusion

We have described a new abstract machine that supports set-valued variables and shown how array manipulating programs can be translated to array-free code for this machine. By compiling array programs for this machine, we are able to discover non-trivial universally quantified loop invariants, simply by analysing the transformed program using off-the-shelf scalar analysers. As an example of how this allows an existing analysis to be lifted to array programs in a straightforward manner, we have extended an uninitialised-variable analysis; Figure 13 showed the usefulness of this approach. The indisputable price for the ease of implementation is a potentially excessive size of the transformed program. However, much array-processing code tends to make simple array traversals and access, and the transformational approach is viable for more than just small programs. Future work includes performing the transformation lazily, to avoid generating unneeded blocks. This should significantly speed up the analysis.

Acknowledgements

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References


Appendix A: Array operations with non-segment variables

Fig. 6(a) assumes that the index variable of every read or write is included in the set of segment bounds. Fig. 14 gives a revised version of \texttt{transform\_body} which handles writes to indices that are not included in the set of segment bounds. When we transform a write to an index in the set of segment bounds (determined by the predicate \texttt{is\_idx}), the transformation is as usual. Otherwise, we emit code to walk through the current set of segments, and apply the write operation to the appropriate one. The dispatch process is similar to the operation of \texttt{split\_transform}, except that all leaves jump back to the continuation of the basic block after the write, rather than continuing under the modified ordering.

\begin{verbatim}
transform\_body(((A[i] = x|stmts], br), \pi, id, ss)
  if is\_idx(i)
    A_i := find\_avar(\pi, A, i)
    transform\_body((stmts, br), \pi, id, ss :=[A_i = x])
  else
    id' := next\_block()
    transform\_body((stmts, br), \pi, id', [])
    dispatch\_write(A[i] = x, \epsilon, \pi, id, ss, id')

dispatch\_write(A[i] = x, sv, [], id, ss, id')
  push\_block(id, ss, Jmp id')

dispatch\_write(A[i] = x, s, [i], id, ss, id')
  id' := next\_block()
  id_ := next\_block()
  s_ := next\_svar(s_)
  s_ := next\_svar(s_)
  if s_ = \epsilon
    push\_block(id, If i < p then id' else id', ss)
  else
    id_ := next\_block()
    push\_block(id, If i < p then id else id', ss)
    push\_block(id_<, Jmp id',
      [(v_1,...,v_A,...,v_k) \in s_<,
       s_ = s_ \cup \{(v_1,...,\pi,...,v_k)\}])
    push\_block(id_> If i = p then id_ else id_>.
      id_ = id_ = id_ = id_ w
    (\ldots,v_A,...) := s_ =
    push\_block(id_<, Jmp id', [v_A = x])
    dispatch\_write(A[i] = x, s_ >, \pi, id_ _, [], id')
\end{verbatim}

Fig. 14. Revised pseudo-code for transforming array writes, allowing for omitted indices. Array reads are transformed similarly.
Analyzing and Compiling Coroutines with Abstract Conjunctive Partial Deduction

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Abstract. We provide an approach to formally analyze the computational behavior of coroutines in Logic Programs and to compile these computations into new programs, not requiring any support for coroutines. The problem was already studied near to 30 years ago, in an analysis and transformation technique called Compiling Control. However, this technique had a strong ad hoc flavor: the completeness of the analysis was not well understood and its symbolic evaluation was also very ad hoc. We show how Abstract Conjunctive Partial Deduction, introduced by Leuschel in 2004, provides an appropriate setting to redefine Compiling Control. Leuschel’s framework is more general than the original formulation, it is provably correct, and it can easily be applied for simple examples. We also show that the Abstract Conjunctive Partial Deduction framework needs some further extension to be able to deal with more complex examples.

1 Introduction

The work reported on in this paper is an initial step in a new project, in which we aim to formally analyze and automatically compile certain types of coroutining computations. Coroutines are a powerful means of supporting complex computation flows. They can be very useful for improving the efficiency of declaratively written programs, in particular for generate-and-test based programs. On the other hand, obtaining a deep understanding of the computation flows underlying the coroutines is notoriously difficult.

In this paper we restrict our attention to pure, definite Logic Programs. In this context, the problem was already studied nearly 30 years ago. [Bruynooghe et al. 1986] and [Bruynooghe et al. 1989] present an analysis and transformation technique for coroutines, called Compiling Control (CC for short). The purpose of the CC transformation is the following: transform a given program, \( P \), into a program \( P' \), so that computation with \( P' \) under the standard selection rule mimics the computation with \( P \) under a non-standard selection rule. In particular, given a coroutining selection rule for a given Logic Program, the transformed program will execute the coroutining if it is evaluated under the standard selection rule of Prolog.
To achieve this, CC consists of two phases: an analysis phase and a synthesis phase. The analysis phase analyzes the computations of a program for a given query pattern and under a (non-standard) selection rule. The query pattern is expressed in terms of a combination of type, mode and aliasing information. The selection rule is instantiation-based, meaning that different choices in atom selection need to be based on different instantiations in these atoms. The analysis results in what is called a “trace tree”, which is a finite upper part of a symbolic execution tree that one can construct for the given query pattern, selection rule and program. In the synthesis phase, a finite number of clauses are generated, so that each clause synthesizes the computation in some branch of the trace tree and such that all computations in the trace tree have been synthesized by some clause. The technique was implemented, formalized and proven correct, under certain fairly technical conditions.

Unfortunately, the CC transformation has a rather ad hoc flavor. It was very hard to show that the analysis phase of the transformation was complete, in the sense that a sufficiently large part of the computation had been analyzed to be able to capture all concrete computations that could possibly occur at run time. Even the very idea of a “symbolic execution” had an ad hoc flavor. It seemed that it should be possible to see this as an instance of a more general framework for analysis of computations.

Fortunately, since the development of CC a number of important advances have been achieved in analysis and transformation:

- General frameworks for abstract interpretation (e.g. [Brunooghe 1991]) were developed. It is clear that abstract interpretation has the potential to provide a better setting for developing the CC analysis. But it still seems different, because abstract interpretation is about analyzing properties that hold during or after a computation, while in CC we are interested in analyzing the computational behavior itself.
- Partial deduction of Logic Programs was developed (e.g. [Gallagher 1986]). Partial deduction seems very similar to CC, but the exact relationship was never identified. When John Lloyd and John Shepherdson formalized the issues of correctness and completeness of partial deduction in [Lloyd and Shepherdson 1991], this provided a new framework for thinking about a complete analysis of a computational behavior and it was clear that some variant of this could improve the CC analysis.
- Conjunctive partial deduction (see [De Schreye et al. 1999]) seems even closer to CC. In an analysis for a CC transformation, one really does not want to split up the conjunctions of atoms into separate ones and then analyze the computations for these atoms separately. It is crucial that one can analyze the computation for certain atoms in conjunction (which is how conjunctive partial deduction generalizes partial deduction), so that their behavior under the non-standard selection rule may be observed.
- Finally, abstract (conjunctive) partial deduction ([Leuschel 2004]) brings all these features together. It provides an extension of (conjunctive) partial
analysing and compiling coroutines

deduction in which the analysis is based on abstract interpretation, rather than on concrete evaluation.

In this paper we will demonstrate – mostly on the basis of examples – that abstract conjunctive partial deduction (ACPD for short) is indeed a suitable framework to redefine CC in such a way that the flaws of the original approach are overcome. We show that for simple problems in the CC context, ACPD can, in principle, produce the transformation automatically. We also show that for more complex CC transformations, ACPD is still not powerful enough. We suggest an extension to ACPD that allows us to solve the problem and illustrate with an example that this extension is very promising.

After the preliminaries, in Section 3 we introduce a fairly refined abstract domain, including type, mode and aliasing information, and we show, by means of an example, how ACPD allows us to analyze a coroutine and compile the transformed program. In Section 4 we propose a more complex example and show why it is out of scope for ACPD. We introduce an additional abstraction in our domain and illustrate that this abstraction solves the problem. This abstraction, however, does not respect the requirements of the formalization of ACPD in Leuschel [2004]. We end with a discussion.

2 Preliminaries

We assume that the reader is familiar with the basics of Logic Programming (Lloyd [1987]). We also assume knowledge of the basics of abstract interpretation (Bruynooghe [1991]) and of partial deduction (Lloyd and Shepherdson [1991]).

In this paper, names of variables will start with a capital. Names of constants will start with a lower case character. Given a program $P$, $Con_p$, $Var_p$, $Fun_p$ and $Pred_p$ respectively denote the sets of all constants, variables, functors and predicate symbols in the language underlying $P$. $Term_p$ will denote the set of all terms constructable from $Con_p$, $Var_p$ and $Fun_p$. $Atom_p$ denotes the set of all atoms which can be constructed from $Pred_p$ and $Term_p$. We will often need to refer to conjunctions of atoms of $Atom_p$ and we denote the set of all such conjunctions as $ConAtom_p$.

We will introduce an abstract domain in the following section. The abstract domain will be based on a set of abstract constant symbols, $ACon_p$. Based on these, there is a corresponding set of abstract terms, $ATerm_p$, which consists of the terms that can be constructed from $ACon_p$ and $Fun_p$. $AAtom_p$ will denote the set of abstract atoms, being the atoms which can be constructed from $ATerm_p$ and $Pred_p$. Finally, $AConAtom_p$ denotes the set of conjunctions of elements of $AAtom_p$.

3 An Example of a CC Transformation, Using ACPD

In this section, we provide the intuitions behind our approach by means of a simple example. We use permutation sort as an illustration. The intention is to transform this program so that calls to perm/2 and ord/1 are interleaved.
Example 1 (Permutation sort).

\[
\begin{align*}
\text{sort}(X,Y) & \leftarrow \text{perm}(X,Y), \text{ord}(Y). \\
\text{perm}([],[]) & . \\
\text{perm}([X|Y],[U|V]) & \leftarrow \text{del}(U,[X|Y],W), \text{perm}(W,V). \\
\text{del}(X,[X|Y],Y) & . \\
\text{del}(X,[Y|U],[Y|V]) & \leftarrow \text{del}(X,U,V). \\
\text{ord}([]) & . \\
\text{ord}([X]) & . \\
\text{ord}([X,Y|Z]) & \leftarrow X \leq Y, \text{ord}([Y|Z]). 
\end{align*}
\]

We now introduce the abstract domain. This domain consists of two types of new constant symbols: \( g \) and \( a_i \), \( i \in \mathbb{N} \). The symbol \( g \) denotes any ground term in the concrete language. The basic intuition for the symbols \( a_i \) is that they are intended to represent variables of the concrete domain. However, as we want the abstract domain to be closed under substitution (if an abstract term denotes some concrete term, then it should also denote all of its instances), an abstract term \( a_i \) will actually represent any term of the concrete language.

The subscript \( i \) in a term \( a_i \) is used to represent aliasing. If an abstract term, abstract atom or abstract conjunction of atoms contains \( a_i \) several times (with the same subscript), the denoted concrete terms, atoms or conjunctions of atoms contain the same term in all positions corresponding to those occupied by \( a_i \).

For instance, the abstract conjunction \( \text{perm}(g,a_1), \text{ord}(a_1) \) denotes the concrete conjunctions \( \{\text{perm}(t_1,t_2), \text{ord}(t_2)|t_1,t_2 \in \text{Term}_p,\ t_1 \text{ is ground}\} \).

In addition to \( g \) and \( a_i \), we will include all concrete constants in the abstract domain, so \( \text{Con}_p \subseteq A\text{Con}_p \). This is not essential for the approach: we could develop a sound and effective ACPD for the CC transformation based on the abstract constants \( g \) and \( a_i \), \( i \in \mathbb{N} \), alone. However, including \( \text{Con}_p \) in \( A\text{Con}_p \) makes the analysis more precise: some redundant paths in the analysis are avoided.

**Definition 1 (Abstract domain).**

The abstract domain consists of:

- \( A\text{Con}_p = \text{Con}_p \cup \{g\} \cup \{a_i|i \in \mathbb{N}\} \).
- \( A\text{Term}_p, A\text{Atom}_p \) and \( A\text{ConAtom}_p \) are defined as the sets of the terms, atoms and conjunctions of atoms constructable from \( A\text{Con}_p, \text{Fun}_p \) and \( \text{Pred}_p \).

Next, we define the semantics of the abstract domain, through a concretization function \( \gamma \). With slight abuse of notation, we use the same symbol \( \gamma \) to denote the concretization functions on \( A\text{Term}_p, A\text{Atom}_p \) and \( A\text{ConAtom}_p \).

In order to formalize the semantics of the aliasing, we need two auxiliary concepts: the subterm selection sequence and the aliasing context.

**Definition 2 (Subterm selection sequence).**

Let \( t \) be a term, atom or conjunction of atoms (either concrete or abstract).

- \( i \in \mathbb{N}_0 \) is a subterm selection sequence for \( t \), if \( t = f(t_1,\ldots,t_n) \) and \( i \leq n \).
- The subterm of \( t \) selected by \( i \) is \( t_i \).
Let $t$ be a subterm selection sequence for $t$, if $t = f(t_1, ..., t_n)$, $i_1 \leq n$, $i_1 \in \mathbb{N}_0$ and $i_2, ..., i_n$ is a subterm selection sequence for $t_{i_1}$. With an inductively defined notation, we denote by $t_{i_1,i_2,...,i_n}$ the subterm of $t_{i_1, ..., i_{n-1}}$ selected by $i_k$, with $1 < k \leq n$. We also refer to $t_{i_1,i_2,...,i_n}$ as the subterm of $t$ selected by $i_1,i_2,...,i_n$.

Note that, in this definition, we assume that a conjunction of atoms $A_1, A_2, ..., A_n$ is denoted as $\wedge(A_1, A_2, ..., A_n)$.

Example 2 (Subterm selection sequence). Let $t = f(g(h(X), 5), f(h(a), Y))$, then $t_{1.1.1} = X$, $t_{2.1.1} = a$.

Definition 3 (Aliasing context).
Let $t$ be an abstract term, atom or conjunction of atoms. The aliasing context of $t$, denoted $AC(t)$, is the finite set of pairs $(sss_1, sss_2)$ of subterm selection sequences of $t$, such that $t_{sss_1} = t_{sss_2} = a_i$ for some $i \in \mathbb{N}$.

Example 3 (Aliasing context).
Let $t = p(f(a_2, g, a_1, a_2, g(h(a_1))))$, then $AC(t) = \{(1.1, 3), (2, 4.1.1)\}$.

Definition 4 (Concretization function).
The concretization function $\gamma : ATerm_p \cup AAtom_p \cup AConAtom_p \rightarrow 2^{Term_p} \cup 2^{Atom_p} \cup 2^{ConAtom_p}$ is defined as:

- $\gamma(c) = \{c\}$, for any $c \in Con_p$
- $\gamma(g) = \{t \in Term_p | t$ is ground$\}$
- $\gamma(a_i) = Term_p$, $i \in \mathbb{N}$
- $\gamma(f(at_1, ..., at_n)) = \{f(t_1, ..., t_n) | t_i \in \gamma(at_i), i = 1...n, and let t denote f(t_1, ..., t_n), then for all (sss_1, sss_2) \in AC(f(at_1, ..., at_n)) : t_{sss_1} = t_{sss_2}\}$

Example 4 (Concretization function).
$\gamma(p(f(a_2, g, a_1, a_2, g(h(a_1)))) = \{p(f(t_1, t_2), t_3, t_1, q(h(t_3))) | t_1, t_3 \in Term_p, t_2$ ground term of $Term_p\}$

The abstract domain introduced above is infinitely large. There are two causes for this. Terms can be nested unboundedly deep, therefore infinitely many different terms exist. In addition, there are infinitely many $a_i, i \in \mathbb{N}$, symbols.

If so desired, the abstract domain can be refined, so that it becomes finite. This is done by using depth-k abstraction and by defining an equivalence relation on $\{a_i | i \in \mathbb{N}\}$. For the purpose of this paper, the infinite size of the abstract domain is not a problem.

Let us return to the permutation sort example. ACPD requires a top-level abstract atom (or conjunction) to start the transformation. Let $sort(g, a_1)$ be this atom. In the context of the $A$-coveredness condition of partial deduction, our initial set $A$ is $\{sort(g, a_1)\}$.

Below, we construct a finite number of finite, abstract partial deduction derivation trees for abstract (conjunctions of) atoms. The construction of these trees assumes an “abstract unification” and an “abstract unfold” operation. Their
formal definitions can be found in Annex [2014]. For now, we only show their effects in abstract partial derivation trees.

Next, we need an “oracle” that decides on the selection rule applied in the abstract derivation trees. This oracle mainly has two functions:

– to decide whether an obtained goal should be unfolded further, or whether it should be kept residual (to be split and added to \(A\)),
– to decide which atom of the current goal should be selected for unfolding.

In fact, we will use a third type of decision that the oracle may make: it may decide to “fully evaluate” a selected atom. This type of decision is not commonly supported in partial deduction. What it means is that we decide not to transform a certain predicate of the original program, but merely keep its original definition in the transformed program. In partial deduction, this can be done by never selecting these atoms, including them in \(A\) and including their original definition in the transformed program.

In our setting, however, we want to know the effect that solving the atom has on the remainder of the goal. Therefore, we will assume that a full abstract interpretation over our abstract domain computes the abstract bindings that solving the atom results in. These are applied to the remainder of the goal. Note that this cannot easily be done in standard partial deduction, as fully evaluating an atom during (concrete) partial deduction may not terminate. In [Vidal 2011], a similar functionality is integrated in a hybrid approach to conjunctive partial deduction.

For now, we simply assume the existence of the oracle. Fig. 1, 2, 3 show the abstract partial derivation trees that ACPD may build for permutation sort and top level \(A = \{\text{sort}(g, a_1)\}\).

---

**Fig. 1:** Abstract tree for \(\text{sort}(g, a_1)\)

\[
\begin{array}{c}
\text{sort}(g, a_1) \\
\text{perm}(g, a_1), \text{ord}(a_1) \\
\mid a_1 = [] \\
\mid \text{ord}([]) \\
\square
\end{array}
\]

**Fig. 2:** Abstract tree for \(\text{perm}(g, a_3), \text{ord}([g|a_3])\)

\[
\begin{array}{c}
\text{perm}(g, a_3), \text{ord}([g|a_3]) \\
\mid a_3 = [] \\
\mid \text{ord}([g]) \\
\mid \text{del}(a_5, g, a_7), \text{perm}(a_7, a_6), \text{ord}([g, a_5|a_6]) \\
\mid a_5 = g, a_7 = g \\
\mid \text{perm}(g, a_6), \text{ord}([g, g|a_6]) \\
\square
\end{array}
\]

In these figures, in each goal, the atom selected for abstract unfolding is underlined. If an atom is underlined twice, this expresses that the atom was selected for full abstract interpretation.

Both unfolding and full abstract evaluation may create bindings. Our abstract unification only collects bindings made on the \(a_i\) terms. Bindings created on \(g\) terms are not relevant.
In the left branch of the tree in Fig. 1 we see the effect of including the concrete constants in the abstract domain. As a result, the binding for \( a_1 \) is \([],\) instead of \( g \). If we had not included \( Con_p \) in \( ACon_p \), then \( ord(g) \) would have required a full analysis, using the three clauses for \( ord/1 \).

A goal with no underlined atom indicates that the oracle selects no atom and decides to keep the conjunction residual. After the construction of the tree in Fig. 1, ACPD adds the abstract conjunction \( perm(g,a_3), ord([g|a_3]) \) to \( A \). ACPD starts a new tree for this atom. This tree is shown in Fig. 2.

The tree is quite similar to the one in Fig. 1. The main difference is that, in the residual leaf, the \( ord \) atom now has a list argument with two \( g \) elements. This pattern does not yet exist in the current \( A \) and is therefore added to \( A \). A third abstract tree is computed for \( perm(g,a_6), ord([g,g|a_6]) \), shown in Fig. 3.

![Fig. 3: Abstract tree for \( perm(g,a_6), ord([g,g|a_6]) \)](image-url)

In Fig. 3, the residual leaf \( perm(g,a_6), ord([g,g|a_6]) \) is a renaming of the conjunction \( perm(g,a_3), ord([g|a_3]) \), which is already contained in \( A \). Therefore, ACPD terminates the analysis, concluding \( A \)-coveredness for \( A = \{sort(g,a_1), \land (perm(g,a_3), ord([g|a_3])), \land (perm(g,a_6), ord([g,g|a_6]))\} \).

In standard (concrete) conjunctive partial deduction, the analysis phase would now be completed. In ACPD, however, we need an additional step. In the abstract partial derivation trees, we have not collected the concrete bindings that unfolding would produce. These are required to generate the resolvents. Therefore, we need an additional step, constructing essentially the same three trees again, but now using concrete terms and concrete unification.

We only show one of these concrete derivation trees in Fig. 4. It corresponds to the tree in Fig. 2. We define the root of a concrete derivation tree corresponding to an abstract tree as follows.

**Definition 5 (Concrete conjunctions in the root).**

Let \(acon\) \(\in A\), then the conjunction in the root of the corresponding concrete tree, denoted as \(c(acon)\), is obtained by replacing any \( g \) or \( a_i \) symbol in \( acon \) by a fresh free variable, ensuring that multiple occurrences of \( a_i \), with the same subscript \( i \), are replaced by identical variables.

When unfolding the concrete tree, every abstract unfolding of the abstract tree is mimicked, using the same clauses, over the concrete domain.
The step of full abstract interpretation of the \textit{del}(a_5, g, a_7) atom in Fig. 2 has no counterpart in Fig. 4. The atom \textit{del}(U, [X_1|X_2], W) is kept residual and the \textit{del/3} clauses are added to the transformed program.

More specifically, using a renaming \( p_1(X, Y, Z) \) for \( \wedge(\text{perm}(X, Y), \text{ord}([Z|Y])) \) and \( p_2(W, V, Z, U) \) for \( \wedge(\text{perm}(W, V), \text{ord}([Z|U])) \), we synthesize the following resolvents from the tree in Fig. 4:

\[
p_1([], [], Z) \leftarrow \text{del}(U, [X_1|X_2], W), p_2(W, V, Z, U).
\]

From the counterparts of the trees in Fig. 1 and Fig. 3, we obtain the following additional resultants:

\[
sort([], []).\]
\[
sort([X_1|X_2], [Y_1|Y_2]) \leftarrow \text{del}(Y_1, [X_1|X_2], Z), p_1(Z, Y_2, Y_1).
\]

\[
p_2(U, V, W, X) \leftarrow W \leq X, p_1(U, V, X).
\]

4 A More Complex Example, Introducing the \textit{multi} Abstraction

In Section 3 we have shown that ACPD is indeed sufficient to formally revisit CC for a simple example. However, for more complex examples, ACPD still lacks expressivity. Consider the following map coloring program.

\textbf{Example 5 (Map coloring).}

\[
safe\text{-}coloring(Ns, Cs) \leftarrow coloring(Cs), safe(Ns, Cs).
\]
\[
coloring([ ]).\]
\[
coloring([C|Cs]) \leftarrow color(C), coloring(Cs).
\]
\[
safe([ ], []).\]
\[
safe([N1, N2|Ns], [C1, C2|Cs]) \leftarrow
\]
\[
\text{allsafe}(N1, C1, [N2|Ns], [C2|Cs]), safe([N3|Ns], [C2|Cs]).
\]
\[
\text{allsafe}(N, C, [ ], []).\]
\[
\text{allsafe}(N1, C1, [N2|Ns], [C2|Cs]) \leftarrow
\]
\[
test(N1, C1, N2, C2), \text{allsafe}(N2, C2, Ns, Cs).
\]
\[
test(N1, C1, N2, C2) \leftarrow edge(N1, N2), C1 \neq C2; \text{nedge}(N1, N2).\]
In addition, we assume a number of facts for edge/2, noedge/2 and color/1, specifying the connections between a number of nations (edge(n₁, n₂) and noedge(n₁, n₂) facts) and a number of colors (color(c) facts). The allsafe/4 predicate tests one nation with respect to all others. The safe/2 predicate tests all nations with respect to one another.

The program is intended to be called with a goal safe_coloring(Ns, Cs), with Ns a ground list of nations without duplicate entries and Cs a list of free variables.

The complete ACPD style analysis is available in Annex [2014]. We only present some relevant parts.

The top level goal for the abstract analysis is safe_coloring(g, a₁), so that the initial set \(A = \{ \text{safe\_coloring}(g, a₁) \}\). A first abstract derivation tree describes the initialization for the computation. It contains two branches leading to an empty goal (success branches) and a third branch with the leaf: \(\land(\text{coloring}(\{a₄|a₅\}), \text{allsafe}(g, g, g, \{a₄|a₅\}), \text{safe}(g, \{a₄|a₅\}))\), which is added to \(A\).

Next, we construct an abstract derivation tree for the latter conjunction. This, again, gives a successful branch, with an empty conjunction in the leaf and a second branch, with the following leaf, which should be added to \(A\): \(\land(\text{coloring}(\{a₆|a₇\}), \text{allsafe}(g, g, g, \{a₆|a₇\}), \text{allsafe}(g, g, g, \{a₆|a₇\}), \text{safe}(g, \{a₆|a₇\}))\).

At this point it becomes clear that an analysis following only the steps shown in Section 3 will not terminate. The two abstract conjunctions, most recently added to \(A\), are identical – up to renaming of \(a_i\)’s – except that the latter conjunction contains two atoms \(\text{allsafe}(g, g, g, \{a_i|a_j\}), \text{allsafe}(g, g, g, \{a_i|a_j\})\)), instead of just one. A further analysis, building additional derivation trees, will result in the construction of continuously growing conjunctions, with continuously increasing numbers of allsafe/4 atoms.

We could solve this by cutting the goal into two smaller conjunctions and adding these to \(A\). However, all these atoms are generators or testers in the coroutine and depend on each other. So, instead of splitting, we extend ACPD. One of the restrictions imposed by ACPD is that for any abstract conjunction of atoms, \(acon \in A\text{ConAtom}_p\), there exists a concrete conjunction, \(con \in \text{ConAtom}_p\), such that: for all \(con_i \in \gamma(acon)\): \(con_i\) is an instance of \(con\). In practice, this means that an abstract conjunction is not allowed to represent a set of concrete conjunctions whose elements have a distinct number of conjuncts. However, in order to solve the problem observed in our example, we need the ability to represent a set of conjunctions, with a growing number of atoms, by an abstract atom.

We propose to introduce a new abstraction into our abstract domain: the multi/1 abstraction.

**Definition 6 (multi/1 abstraction).**
We extend \(A\text{ConAtom}_p\) with a new construct \(\text{multi}(A)\), where \(A \in A\text{Atom}_p\), such that: \(\gamma(\text{multi}(A)) = \{ \gamma(\land_{i=1\ldots n} A) | n \in \mathbb{N}_0 \}\).

**Example 6 (multi/1 abstraction).**
The abstract conjunction \(\text{multi}(\text{allsafe}(g, g, g, \{a₆|a₇\}))\) denotes all finite concrete conjunctions \(AS_1 \land \ldots \land AS_n, n \in \mathbb{N}_0, AS_i \in \gamma(\text{allsafe}(g, g, g, \{a₆|a₇\}))\) for each \(i = 1, n\).
We now define two new operations on our abstract domain: abstract unfolding of \textit{multi/1} and abstract generalization with \textit{multi/1}.

**Definition 7 (Abstract unfolding of \textit{multi/1}).**
Let \( A \in AAtom_p \). Abstract unfolding of \( \text{multi}(A) \) is defined as standard unfolding using the clauses:
\[
\text{multi}(X) \leftarrow X.
\]
\[
\text{multi}(X) \leftarrow X, \text{multi}(X).
\]

The definition states that unfolding an abstract \( \text{multi}(A) \) conjunction makes a case split: the case in which there is only one atom \( A \) in the abstract conjunction and the case in which there are more than one.

**Definition 8 (**\textit{multi/1} generalization).**
Let \( A \in AAtom_p \) and \( n \in \mathbb{N}_0 \). Abstract generalization with \textit{multi/1} is defined as the function \( \text{Gen} : AConAtom_p \rightarrow AConAtom_p \):
\[
\text{Gen}(\bigland_{i=1}^n A) = \text{multi}(A)
\]
\[
\text{Gen}((\bigland_{i=1}^n A) \land \text{multi}(A)) = \text{multi}(A)
\]
\[
\text{Gen}((\text{multi}(A) \wedge (\bigland_{i=1,...,n} A)) = \text{multi}(A)
\]

We illustrate these operations in our running example. Observing the growing number of \( \text{allsafe}(g,g,g,[a_6|a_7]) \) atoms in our last conjunction (w.r.t. the conjunction already present in \( A \)), we perform the generalization: \( \text{Gen}((\bigland(\text{allsafe}(g,g,g,[a_6|a_7]), \text{allsafe}(g,g,g,[a_6|a_7]))) = \text{multi}(\text{allsafe}(g,g,g,[a_6|a_7]))
\]

We replace the conjunction \( \bigland(\text{coloring([a_4|a_5]), allsafe}(g,g,g,[a_4|a_5]), \text{safe}(g,[a_4|a_5])) \) from \( A \) by:
\[
\bigland(\text{coloring([a_6|a_7]), allsafe}(g,g,g,[a_6|a_7]), \text{safe}(g,[a_6|a_7]))
\]

Then we construct a new abstract derivation tree for this conjunction, including – among other – an abstract unfold of \textit{multi/1} and abstract generalizations with \textit{multi/1}. In Fig. 5, 6, 7 we show this abstract tree.

After abstract unfolding of \textit{coloring([a_1|a_2])} and full abstract evaluation of \textit{color(a_1)}, the tree contains an abstract unfolding of \textit{multi(allsafe}(g,g,g,[g|a_2])).

Fig. 6 contains the base case for this unfolding. After an unfold of \textit{allsafe/4} and full solve of \textit{test/4}, unfolding of \textit{safe}(g,[g|a_2]), using the clause \textit{safe}([N],[C]) \leftarrow \ldots, eventually leads to the empty goal. Unfolding \textit{safe}(g,[g|a_2]) with the recursive clause produces a goal with two \textit{allsafe}(g,g,g,[a_3|a_1]) atoms. We perform generalization with \textit{multi/1}.

In Fig. 7 the non-base case for the abstract unfolding of \textit{multi/1} unfolds an \textit{allsafe}(g,g,g,[g|a_2]) and fully solves a \textit{test}(g,g,g,g). Finally, we perform a generalization with \textit{multi/1} for the \textit{allsafe}(g,g,g,a_2) atom, reaching the new conjunction of \( A \):
\[
\bigland(\text{coloring(a_2)}, \text{multi(allsafe}(g,g,g,a_2)), \text{multi}(\text{allsafe}(g,g,g,[g|a_2])), \text{safe}(g,[g|a_2]))
\]

Eventually, the analysis ends up with a final set \( A \):
\[
\{ \bigland(\text{safe\_coloring}(g,a_1)), \text{multi(allsafe}(g,g,[a_1])), \bigland(\text{coloring}(a_1[a_2]), \text{multi(allsafe}(g,g,g,a_1[a_2])), \text{safe}(g,[a_1[a_2])), \bigland(\text{coloring}(a_2), \text{multi(allsafe}(g,g,a_2)), \text{multi}(\text{allsafe}(g,g,[g|a_2])), \text{safe}(g,[g|a_2])) \}
\]
Analysing and Compiling Coroutines

**Fig. 5:** A recurrent pattern in which a *multi* abstraction is unfolded by case analysis

\[
\begin{array}{c}
\text{coloring(}a_1|a_2\text{)} , \text{multi(}\text{allsafe}(g,g,g,|a_1|a_2)\text{)} , \text{safe}(g,|a_1|a_2) \\
\text{coloring(}a_1\text{)} , \text{coloring}(a_2) , \text{multi(}\text{allsafe}(g,g,g,|a_1|a_2)\text{)} , \text{safe}(g,|a_1|a_2) \\
\text{unfold multi} \\
an_1 = g \\
\text{coloring}(a_2) , \text{multi(}\text{allsafe}(g,g,|g|a_2)\text{)} , \text{safe}(g,|g|a_2) \\
\end{array}
\]

One **allsafe/4**: see Fig. 6

Multiple **allsafe/4**: see Fig. 7

**Fig. 6:** First case from Fig. 5

All non-empty leaves in the abstract derivation trees for these atoms are (renamings of) elements of \(A\). This shows \(A\)-coveredness and the abstract phase of the analysis terminates.

Similar to what was observed for permutation sort in Section 3, we still need an extra analysis to collect the concrete bindings, so that the resultants can be generated. Special care is required for the *multi/1* abstraction. There are three issues: how to represent *multi/1* in the concrete domain, how to deal with the concrete counterparts of abstract generalization with *multi/1* and abstract unfolding of *multi/1*.

Definition 5, in Section 3, defined the concrete counterparts of the conjunctions in \(A\). We extend it to *multi*(\(A\)):

**Definition 9 (Concrete conjunction for *multi*(\(A\))).**
Let \(A \in \text{AAAtom}_p\), then \(c(\text{multi}(A)) = \text{multi}([c(A)|T])\), with \(T\) a fresh variable.

**Example 7.** \(c(\text{multi(}\text{allsafe}(g,g,g,a_2)\text{)}) = \text{multi(}\text{allsafe}(X,Y,U,V)|T)\)

For the abstract generalization with *multi/1*, we define the concrete counterpart as follows.
Definition 10 (Concrete generalization). Let $A \in AAtom$.

- If the abstract generalization with $\text{multi}/1$ is of the type $\text{Gen}(\bigwedge_{i=1}^{n} A) = \text{multi}(A)$, then the corresponding node in the concrete derivation contains $c((\bigwedge_{i=1}^{n} A))$, where $c(A)$ is a list of at least one $c(A)$. The concrete generalization is defined as $\text{ConGen}(c((\bigwedge_{i=1}^{n} A))) = \text{multi}(c([A, ..., A]))$, with $n$ members in the list.
- If the abstract generalization with $\text{multi}/1$ is of the type $\text{Gen}((\bigwedge_{i=1}^{n} A) \land \text{multi}(A)) = \text{multi}(A)$, then the corresponding node in the concrete derivation contains $c((\bigwedge_{i=1}^{n} A) \land \text{multi}(\text{List}))$, where $\text{List}$ is a list of at least one $c(A)$. The concrete generalization is defined as $\text{ConGen}(c((\bigwedge_{i=1}^{n} A) \land \text{multi}(\text{List}))) = \text{multi}(c([A, ..., A][\text{List}]))$ with $n$ new members added to $\text{List}$.
- The third case, $\text{Gen}(\text{multi}(A) \land (\bigwedge_{i=1}^{n} A)) = \text{multi}(A)$, is treated identically to the previous one.

Example 8 (Concrete generalization). Let $\text{allsafe}(X_1, Y_1, [U_1|V_1], [W|Z]), \text{allsafe}(X_2, Y_2, [U_2|V_2], [W|Z])$ occur in a concrete conjunction in a concrete derivation tree, where abstract generalization with $\text{multi}/1$ is performed on the corresponding abstract conjunction. Then, as a next step in the concrete derivation tree, this conjunction is replaced by $\text{multi}(\text{allsafe}(X_1, Y_1, [U_1|V_1], [W|Z]), \text{allsafe}(X_2, Y_2, [U_2|V_2], [W|Z]))$.

Note that this “generalization” actually does not generalize anything. It only brings the information in a form that can be generalized.

The actual generalization happens implicitly in the move to the construction of the next concrete derivation tree. If our conjunction is a leaf of the concrete derivation tree, then the corresponding abstract conjunction is added to the set $A$. Let $\wedge(\text{coloring}([a_1|a_2]), \text{multi}(\text{allsafe}(g, g, [a_1|a_2])), \text{safe}(g, [a_1|a_2]))$, for instance, be the corresponding abstract conjunction that is added to $A$. Then, a new concrete tree is built for a concrete atom corresponding to this abstract one.

In this example, the root of that concrete tree is:

$$\wedge(\text{coloring}([W|Z]), \text{multi}(\text{allsafe}(X, Y, V, [W|Z])|\text{Tail})), \text{safe}(U, [W|Z]))$$

Finally, we still need to define the counterpart of abstract unfold of $\text{multi}/1$ in the concrete tree. To do this, we add the following definition of $\text{multi}/1$ to the original program $P$:
\begin{verbatim}
multi([X]) ← X.
multi([X|Y]) ← X, multi(Y).
\end{verbatim}

It should be clear that mere concrete unfolding of concrete \texttt{multi/1} atoms with the above definition for \texttt{multi/1} gives us the desired counterpart of the case split performed in abstract unfold of \texttt{multi/1}.

With the concepts above, we construct a concrete derivation tree, mimicking the steps in the abstract derivation tree – but over the concrete domain – for every conjunction in the set \( \mathcal{A} \). Collecting all the resultants from these concrete trees, we get the transformed program. A working Prolog program can be found in Annex [2014].

5 Discussion

In this paper, we have presented an approach to formally analyze the computations, for Logic Programs, performed under coroutining selection rules, and to compile such computations into new programs. On the basis of an example, we have shown that simple coroutines, in which the execution of a single, atomic generator is interleaved with a single, atomic tester, can be successfully analyzed and compiled within the framework of ACPD (Leuschel [2004]). These “simple” coroutines essentially correspond to the \textit{strongly regular} logic programs of Vidal [2011].

To achieve this, we defined an expressive abstract domain, capturing modes, types and aliasing. In the paper, we have focused on the intuitions, more than on the full formalization, as space restrictions would not allow both. However, we have developed the formal definitions for the ordering on the abstract domain, abstract unification, abstract unfold and others.

Because the approach – for simple coroutines – fits fully within the ACDP framework, it inherits the correctness results from ACPD. In particular, \( \mathcal{A} \)-closedness and independence guarantee the completeness and correctness of the analysis. In addition, the transformation preserves all computed answers (in both directions) and finite failure of the transformed program implies finite failure of the original.

We have proposed an extension to our abstract domain: the \texttt{multi/1}-abstraction. A \texttt{multi/1} atom can represent (sets of) conjunctions of one or more concrete atoms. We have defined abstract unfold and abstract generalisation operations for this abstraction. We have shown, in an example, that this abstraction and these operations allow us to extend ACPD, enabling it to perform a complete analysis, and to compile the more complex coroutines.

On a more general level, our work provides a new, rational reconstruction of the CC-transformation (Bruynooghe et al. [1986]), avoiding ad hoc features of the CC approach. In addition, the work presents a new application for ACPD.

As a rule, coroutining improves the efficiency of declarative programs by testing partial solutions as quickly as possible. In addition, a program may become more flexible when the transformation is applied. Recall that the original is meant to be called with a ground list of nations and a list of free variables. The
transformed program can be run in the same way, but the top-level predicate can also be called with a ground list of nations and a free variable. This is because SLD resolution sends the original program down an infinite branch of the search tree. The transformed program checks results earlier and, as a result, infers that both top-level arguments must be lists of the same size. In this scenario, compiling control transforms an infinite computation into a finite one.

The CC-transformation raised challenges for a number of researchers and a range of compediting transformation and synthesis techniques. A first reformulation of the CC-transformation was proposed in the context of the “programs-as-proofs” paradigm, in Wiggins [1990]. It was shown that CC-transformations, to a limited extent, could be formalized in a proof-theoretic program synthesis context.

In Boulanger et al. [1993], CC-transformation was revisited on the basis of a combination of abstract interpretation and constraint processing. This improved the formalization of the technique, but it did not clarify the relation with partial deduction.

The seminal survey paper on Unfold/Fold transformation, Pettorossi and Proietti [1994], showed that basic CC-transformations are well in the scope of Unfold/Fold transformation. In later works (e.g. Pettorossi and Proietti [2002]), the same authors introduced list-introduction into the Unfold/Fold framework, whose function is very similar to that of the multi/1 abstraction in our approach. Also related to our work are Puebla et al. [1997], providing alternative transformations to improve the efficiency of dynamic scheduling, and Vidal [2011] and Vidal [2012], which also provide a hybrid form of partial deduction, combining abstract and concrete levels.

There are a number of issues that are open for future research. First, we aim to investigate the generality of the multi/1 abstraction. Although it seems to work well in a number of examples, we will study more complex ones. We also want to revisit the ACPD framework, in order to extend it to the new abstraction we aim to support. This will involve a new formalization of ACPD, capable of supporting analysis and compilation of coroutines in full generality. This will also formally establish the correctness results for the more general cases, such as the one presented in Section 4. Obviously, we also want to have a full implementation of these concepts and to show that the analysis and compilation can be fully automated.

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References


Confluence Modulo Equivalence in Constraint Handling Rules

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Abstract. Previous results on confluence for Constraint Handling Rules, CHR, are generalized to take into account user-defined state equivalence relations. This allows a much larger class of programs to enjoy the advantages of confluence, which include various optimization techniques and simplified correctness proofs. A new operational semantics for CHR is introduced that reduces notational overhead significantly and allows to consider confluence for programs with extra-logical and incomplete built-in predicates. Proofs of confluence are demonstrated for programs with redundant data representation, e.g., sets-as-lists, for dynamic programming algorithms with pruning as well as a Union-Find program, which are not covered by previous confluence notions for CHR.

1 Introduction

A rewrite system is confluent if all derivations from a common initial state end in the same final state. Confluence, like termination, is often a desirable property, and proof of confluence is a typical ingredient of a correctness proof. For a programming language based on rewriting such as Constraint Handling Rules, CHR, it ensures correctness of parallel implementations and application order optimizations.

Previous studies of confluence for CHR programs are based on Newman’s lemma. This lemma concerns confluence defined in terms of alternative derivations ending in the exact same state, which excludes a large class of interesting CHR programs. However, the literature on confluence in general rewriting systems has, since the early 1970s, offered a more general notion of confluence modulo an equivalence relation; this defines that alternative derivations only need to end in states that are equivalent with respect to some equivalence relation (and not necessarily identical). In this paper, we show how confluence modulo equivalence can be applied in a CHR context, and we demonstrate interesting programs covered by this notion that are not confluent in any previous

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versions. The use of redundant data representations is one example of what becomes within reach of confluence, and programs that search for one best among multitudes of alternative solutions is another.

**Example 1.** The following CHR program, consisting of a single rule, collects a number of separate items into a (multi-) set represented as a list of items.

\[
\text{set}(L), \text{item}(A) \iff \text{set}([A|L]).
\]

This rule will apply repeatedly, replacing constraints matched by the left hand side by those indicated to the right. The query

\[?- \text{item}(a), \text{item}(b), \text{set}([]).\]

may lead to two different final states, \{set([a,b])\} and \{set([b,a])\}, both representing the same set. This can be formalized by a state equivalence relation \(\approx\) that implies \{set(L)\} \(\approx\) \{set(L')\}, whenever \(L\) is a permutation of \(L'\). The program is not confluent when identical end states are required, but it will be shown to be confluent modulo \(\approx\).

Our generalization is based on a new operational semantics that permits extra-logical and incomplete predicates (e.g., Prolog's \text{var}/2 and \text{is}/2), which is out of the scope of previous approaches. It also leads to a noticeable reduction of notational overhead due to a simpler structure of states.

Section 2 reviews previous work on confluence, in general and for CHR. Sections 3 and 4 give preliminaries and our operational semantics. Section 5 considers how to prove confluence modulo equivalence for CHR. Section 6 shows confluence modulo equivalence for a version in CHR of the Viterbi algorithm; it represents a wider class of dynamic programming algorithms with pruning, also outside the scope of earlier proposals. Section 7 shows confluence modulo equivalence for the Union-Find algorithm, which has become a standard test case for confluence in CHR; it is not confluent in any previously proposed way (except with construed side-conditions). Section 8 comments on related work in more details, and the final section provides a summary and a conclusion.

### 2 Background

A binary relation \(\rightarrow\) on a set \(A\) is a subset of \(A \times A\), where \(x \rightarrow y\) denotes membership of \(\rightarrow\). A rewrite system is a pair \(<A, \rightarrow>\); it is terminating if there is no infinite chain \(a_0 \rightarrow a_1 \rightarrow \cdots\). The reflexive transitive closure of \(\rightarrow\) is denoted \(\rightarrow^*\). The inverse relation \(\leftarrow\) is defined by \(\{(y, x) \mid x \rightarrow y\}\). An equivalence (relation) \(\approx\) is a binary relation on \(A\) that is reflexive, transitive and symmetric.

A rewrite system \(<A, \rightarrow>\) is confluent if and only if \(y \leftarrow x \rightarrow y' \Rightarrow \exists z. y \rightarrow z \leftarrow z'\), and is locally confluent if and only if \(y \leftarrow x \rightarrow y' \Rightarrow \exists z. y \rightarrow z \leftarrow z'\). In 1942, Newman showed his fundamental lemma [1]: A terminating rewrite system is confluent if and only if it is locally confluent. An elegant proof of Newman’s lemma was provided by Huet [2] in 1980.

The more general notion of confluence modulo equivalence was introduced in 1972 by Aho et al. [3] in the context of the Church-Rosser property.
Definition 1 (Confluence modulo equivalence). A relation $\rightarrow$ is confluent modulo an equivalence $\approx$ if and only if

$$\forall x, y, x', y'. \ x' \xrightarrow{\approx} x \approx y \xrightarrow{\approx} y' \Rightarrow \exists z, z'. \ x' \xrightarrow{\approx} z \approx z' \xleftarrow{\approx} y'$$

This shown as a diagram in Fig. 1a. Sethi [4] showed in 1974 that confluence modulo equivalence for a bounded rewrite system is equivalent to the following properties, $\alpha$ and $\beta$; see also Fig. 1b.

Definition 2 ($\alpha$ & $\beta$). A relation $\rightarrow$ has the $\alpha$ property and the $\beta$ property if and only if it satisfy the $\alpha$ condition and the $\beta$ condition, respectively:

$$\alpha: \ \forall x, y, y'. \ y' \xleftarrow{x} x \rightarrow y \Rightarrow \exists z, z'. \ y' \xrightarrow{z} z \approx z' \xleftarrow{y} y$$

$$\beta: \ \forall x, x', y. \ x' \approx x \rightarrow y \Rightarrow \exists z, z'. \ x' \xrightarrow{z} z \approx z' \xleftarrow{y} y$$

In 1980, Huet [2] generalized this result to any terminating system.

Definition 3 (Local confl. mod. equivalence). A rewrite system is locally confluent modulo an equivalence $\approx$ if and only if it has the $\alpha$ and $\beta$ properties.

Theorem 1. Let $\rightarrow$ be a terminating relation. For any equivalence $\approx$, $\rightarrow$ is confluent modulo $\approx$ if and only if $\rightarrow$ is locally confluent modulo $\approx$.

Fig. 1: Diagrams for the fundamental notions. A dotted arrow (single wave line) indicates an inferred step (inferred equivalence).

The known results on confluence for CHR are based on Newman’s lemma. Abdennadher et al [5] in 1996 seem to be the first to consider this, and they could show that confluence (without equivalence) for CHR is decidable and can be checked by examining a finite set of states formed by a combination of heads of rules. A refinement, called observational confluence was introduced in 2007 by Duck et al [6], in which only states that satisfy a given invariant are considered.

3 Preliminaries

We assume standard notions of first-order logic such as predicates, atoms and terms. For any expression $E$, $\text{vars}(E)$ refers to the set of variables that occurs
in \( E \). A substitution is a mapping from a finite set of variables to terms, which also may be viewed as a set of first-order equations. For substitution \( \sigma \) and expression \( E \), \( E\sigma \) (or \( E \cdot \sigma \)) denotes the expression that arises when \( \sigma \) is applied to \( E \); composition of two substitutions \( \sigma, \tau \) is denoted \( \sigma \circ \tau \). Special substitutions \textit{failure}, \textit{error} are assumed, the first one representing falsity and the second runtime errors.

Two disjoint sets of \textit{(user) constraints} and \textit{built-in} predicates are assumed. For the built-ins, we use a semantics that is more in line with implemented CHR systems than previous approaches, and which allows not only logical built-ins but also extra-logical devices such as Prolog’s \texttt{var/1} and incomplete ones such as \texttt{is/2}. While \cite{5,6,7} collect built-ins in a separate store and determine their satisfiability by a magic solver that mirrors a first-order semantics, we execute a built-in right away, which means that it serves as a test, possible giving rise to a substitution that is applied to form the subsequent state.

An evaluation procedure \( \text{Exe} \) for built-ins \( b \) is assumed, such that \( \text{Exe} (b) \) is either a (possibly identity) substitution to a subset of \( \text{vars}(b) \) or one of \textit{failure}, \textit{error}. It extends to sequences of built-ins as follows.

\[
\text{Exe}(b_1, b_2) = \begin{cases} 
\text{Exe}(b_1) & \text{when } \text{Exe}(b_1) \in \{\text{failure}, \text{error}\}, \\
\text{Exe}(b_2 \cdot \text{Exe}(b_1)) & \text{when otherwise } \text{Exe}(b_2 \cdot \text{Exe}(b_1)) \\
\text{Exe}(b_1) \circ \text{Exe}(b_2 \cdot \text{Exe}(b_1)) & \text{otherwise}
\end{cases}
\]

A subset of built-in predicates are the \textit{logical} ones, whose meaning is given by a first-order theory \( B \). For a logical atom \( b \) with \( \text{Exe}(b) \neq \text{error} \), the following conditions must hold.

- Partial correctness: \( B \models \forall_{\text{vars}(b)} (b \leftrightarrow \exists_{\text{vars}(\text{Exe}(b)) \setminus \text{vars}(b)} \text{Exe}(b)) \).
- Instantiation monotonicity: \( \text{Exe}(b \cdot \sigma) \neq \text{error} \) for all substitutions \( \sigma \).

A logical predicate \( p \) is \textit{complete} whenever, for any \( p \) atom \( b \) that \( \text{Exe}(b) \neq \text{error} \); later we define completeness with respect to a state invariant. Any built-in predicate which is not logical is called \textit{extra-logical}. The following predicates are examples of built-ins; \( \epsilon \) is the empty substitution.

1. \( \text{Exe}(t = t') = \sigma \) where \( \sigma \) is a most general unifier of \( t \) and \( t' \); if no such unifier exists, the result is \textit{failure}.
2. \( \text{Exe}(\text{true}) \) is \( \epsilon \).
3. \( \text{Exe}(\text{fail}) \) is \textit{failure}.
4. \( \text{Exe}(t \text{ is } t') = \text{Exe}(t = v) \) whenever \( t' \) is a ground term that can be interpreted as an arithmetic expression \( e \) with the value \( v \); if no such \( e \) exists, the result is \textit{error}.
5. \( \text{Exe}(\text{var}(t)) = \epsilon \) if \( t \) is a variable and \textit{failure} otherwise.
6. \( \text{Exe}(\text{ground}(t)) = \epsilon \) when \( t \) is ground and \textit{failure} otherwise.
7. \( \text{Exe}(t \approx t') = \epsilon \) when \( t \) and \( t' \) are identical and \textit{failure} otherwise.
8. \( \text{Exe}(t \neq t') = \epsilon \) when \( t \) and \( t' \) are non-unifiable and \textit{failure} otherwise.
The first three predicates are logical and complete; “is” is logical but not complete without an invariant that grounds its second arguments (considered later). The remaining ones are extra-logical.

The practice in previous semantics [5,6,7] of conjoining built-ins and testing them by satisfiability leads to ignorance of runtime errors and incompleteness.

To represent the propagation history, we introduce indices: An indexed set $S$ is a set of items of the form $x:i$ where $i$ belongs to some index set and each such $i$ is unique in $S$. When clear from context, we may identify an indexed set $S$ with its cleaned version $\{x \mid x:i \in S\}$. Similarly, the item $x$ may identify the indexed version $x:i$. We extend this to any structure built from indexed items.

## 4 Constraint Handling Rules

We define an abstract syntax of CHR together with an operational semantics suitable for considering confluence. We use the generalized simpagation form as a common representation for the rules of CHR. Guards may unify variables that occur in rule bodies, but not variables that occur in the matched constraints. In accordance with the standard behaviour of implemented CHR systems, failure and runtime errors are treated the same way in the evaluation of a guard, but distinguished when occurring in a query or rule body, cf. definition 7, below.

**Definition 4.** A rule $r$ is of the form

$$H_1 \setminus H_2 \Leftarrow g \mid C,$$

where $H_1$ and $H_2$ are sequences of constraints, forming the head of the rule, $g$ is a guard being a sequence of built-ins, and $C$ is a sequence of constraints and built-ins called the body of $r$. Any of $H_1$ and $H_2$, but not both, may be empty.

A program is a finite set of rules.

For any fresh variant of rule $r$ with notation as above, an application instance $r''$ is given as follows.

1. Let $r'$ be a structure of the form

$$H_1 \setminus H_2 \Leftarrow C\tau\sigma,$$

where $\tau$ is a substitution for the variables of $H_1$, $H_2$, $\text{Exe}(g\tau) = \sigma \notin \{\text{failure}, \text{error}\}$, and it holds that $(H_1 \setminus H_2)\tau = (H_1 \setminus H_2)\tau\sigma$,

2. $r''$ is a copy of $r'$ in which each atom in its head and body is given a unique index, where the indices used for the body are new and unused.

The application record for such an application instance is of the form

$$r @ i_1, \ldots, i_n$$

where $i_1, \ldots, i_n$ is the sequence of indices of $H_1, H_2$ in the order they occur.

Guards are removed from application instances as they are a priori satisfied.

A rule is a simplification when $H_1$ is empty, a propagation when $H_2$ is empty; in both cases, the backslash is left out, and for a propagation, the arrow symbol is written $==$>. Any other rule is a simpagation. The following definition will become useful later on when we consider confluence.
Definition 5. Consider two application instances \( r_i = (A_i \setminus B_i, \leq C_i), i = 1, 2 \). We say that \( r_1 \) is blocking \( r_2 \) whenever \( B_1 \cap (A_2 \cup B_2) \neq \emptyset \).

For this to be the case, \( r_1 \) must be a simplification or simpagation. Intuitively, it means that if \( r_1 \) has been applied to a state, it is not possible subsequently to apply \( r_2 \). In the following definition of execution states for CHR, irrelevant details of the state representation are abstracted away using principles of [8]. To keep notation consistent with Section 2, we use letters such as \( x, y \), etc. for states.

Definition 6. A state representation is a pair \( \langle S, T \rangle \), where

- \( S \) is a finite, indexed set of atoms called the constraint store,
- \( T \) is a set of application records called the propagation history.

Two state representations \( S_1 \) and \( S_2 \) are isomorphic, denoted \( S_1 \equiv S_2 \) whenever one can be derived from the other by a renaming of variables and a consistent replacement of indices (i.e., by a 1-1 mapping). When \( \Sigma \) is the set of all state representations, a state is an element of \( \Sigma/\equiv \cup \{ \text{failure}, \text{error} \} \), i.e., an equivalence class in \( \Sigma \) induced by \( \equiv \) or one of two special states; applying the failure (error) substitution to a state yields the failure (error) state. To indicate a given state, we may for simplicity mention one of its representations.

A query \( q \) is a conjunction of constraints, which is also identified with an initial state \( \langle q', \emptyset \rangle \) where \( q' \) is an indexed version of \( q \).

To make statements about, say, two states \( x, y \) and an instance of a rule \( r \), we may do so mentioning state representatives \( x', y' \) and application instance \( r' \) having recurring indices.

In contrast to [5,6,7], we have excluded global variables as they are easy to simulate: A query \( q(X) \) is extended to \( \text{global}(X', X), q(X) \), where \( \text{global}/2 \) is a new constraint predicate; \( X' \) is a constant that serves as a name of the variable. The value \( \text{val} \) for \( X \) is found in the final state in the unique constraint \( \text{global}(X', \text{val}) \). References [5,6,7] use a state component for constraints waiting to be processed, plus a separate derivation step to introduce them into the constraint store. This is unnecessary as the derivations made under either premises are basically the same. Our derivation relation is defined as follows; here and in the rest of this paper, \( \cup \) denotes union of disjoint sets.

Definition 7. A derivation step \( \rightarrow \) from one state to another can be of two types: by rule \( \rightarrow^r \) or by built-in \( \rightarrow^b \), defined as follows.

Apply: \( \langle S \cup H_1 \cup H_2, T \rangle \rightarrow^r \langle S \cup H_1 \cup C, T' \rangle \)

whenever there is an application instance \( r \) of the form \( H_1 \setminus H_2 \leq C \) with \( \text{applied}(r) \notin T \), and \( T' \) is derived from \( T \) by 1) removing any application record having an index in \( H_2 \) and 2) adding \( \text{applied}(r) \)

in case \( r \) is a propagation.

Built-in: \( \langle \{ b \} \cup S, T \rangle \rightarrow^b \langle S, T \rangle \cdot \text{Exe}(b) \).
A state $z$ is final for query $q$, when $q \rightarrow^* z$ and no step is possible from $z$.

The removal of certain application records in Apply steps means to keep only those records that are essential for preventing repeated application of the same rule to the same constraints (identified by their indices).

As noticed by [6], an invariant can make more programs confluent as unusual states, that never appears in practice, are avoided. An invariant may also make it easier to characterize an equivalence relation for states.

**Definition 8.** An invariant is a property $I(\cdot)$ which may or may not hold for a state, such that $I(x) \land x \rightarrow y \Rightarrow I(y)$. An $I$-state is a state $x$ with $I(x)$, and an $I$-derivation is one starting from an $I$-state. A program is $I$-terminating whenever all $I$-derivations are terminating. A set of allowed queries $Q$ may be specified, giving rise to an invariant reachable $Q(x) \Leftrightarrow \exists q \in Q: q \rightarrow^* x$.

A (state) equivalence is an equivalence relation $\approx$ on the set of $I$-states.

Theorem 1 applies specifically for CHR programs equipped with invariant $I$ and equivalence relation $\approx$. When $\approx$ is identity, it coincides with a theorem of [6] for observable confluence. If, furthermore, $I \Leftrightarrow true$, we obtain the classical confluence results for CHR [9].

The following definition is useful when considering confluence for programs that use Prolog built-ins such as “is/2”.

**Definition 9.** A logical predicate $p$ is complete with respect to invariant $I$ whenever, for any $p$ atom $b$ in some $I$-state that $\text{Exe}(b) \neq \text{error}$.

As promised earlier, “is/2” is complete with respect to an invariant that guarantees groundness of the second argument of any call to “is/2”.

**Example 2.** Our semantics permits CHR programs that define constraints such as Prolog’s dif/2 constraint and a safer version of is/2.

\[
\begin{align*}
dif(X,Y) & \iff X=Y \mid \text{fail}. \\
dif(X,Y) & \iff X\neq Y \mid \text{true}. \\
X \text{ safer_is } Y & \iff \text{ground}(Y) \mid X \text{ is } Y.
\end{align*}
\]

5 Proving Confluence Modulo Equivalence for CHR

We consider here ways to prove the local confluence properties $\alpha$ and $\beta$ from which confluence modulo equivalence may follow, cf. Theorem 1. The corners in the following definition generalize the critical pairs of [5]. For ease of usage, we combine the common ancestor states with the pairs, thus the notion of corners corresponding to the “given parts” of diagrams for the $\alpha$ and $\beta$ properties, cf. Fig. 1a. The definitions below assume a given $I$-terminating program with invariant $I$ and state equivalence $\approx$. Two states $x$ and $x'$ are joinable modulo $\approx$ whenever there exist states $z$ and $z'$ such that $x \rightarrow^* z \approx z' \leftarrow^* x'$. 
Definition 10. An α-corner consists of I-states \( x, y \) and \( y' \) with \( y \neq y' \) and two derivation steps such that \( y \xrightarrow{\delta} x \xrightarrow{\delta} y' \). An α-corner is joinable modulo \( \approx \) whenever \( y \) and \( y' \) are joinable modulo \( \approx \).

A β-corner consists of I-states \( x, x' \) and \( y \) with \( x \neq x' \) and a derivation step such that \( x' \approx x \xrightarrow{\gamma} y \). A β-corner is joinable modulo \( \approx \) whenever \( x' \) and \( y \) are joinable modulo \( \approx \).

Some corners are critical, meaning that their satisfaction of the α or β property is not trivial.

Definition 11. An α-corner \( y \xrightarrow{\delta} x \xrightarrow{\delta} y' \) or β-corner \( x' \approx x \xrightarrow{\gamma} y \) is critical if one of the following properties holds.

\( \alpha_1 \): \( \gamma \) and \( \delta \) are application instances where \( \gamma \) blocks \( \delta \) (Def. 5).

\( \alpha_2 \): \( \gamma \) is an application instance of a rule whose guard \( g \) contains an extra-logical built-in, and \( \delta \) is a built-in with \( \text{vars}(g) \cap \text{vars}(\delta) \neq \emptyset \).

\( \alpha_3 \): \( \gamma \) and \( \delta \) are built-ins with \( \gamma \) being extra-logical or not complete wrt. \( I \), and \( \text{vars}(\gamma) \cap \text{vars}(\delta) \neq \emptyset \).

\( \beta \): there exists no state \( y' \) and single derivation step of \( x' \) such that \( x' \xrightarrow{\delta} y' \approx y \).

Critical β-corners are motivated by the experience that often the \( \delta \) step can be formed trivially by applying the same rule or built-in of \( \gamma \) in an analogous way to the state \( x' \). By inspection and Theorem 1, we get the following.

Lemma 1. Any non-critical corner is joinable modulo \( \approx \).

Theorem 2. A terminating program is confluent modulo \( \approx \) if and only if all its critical corners are joinable modulo \( \approx \).

Without invariant, equivalence and extra-logicals, the only critical corners are of type \( \alpha_1 \); here [5] has shown that joinability of a finite set of minimal critical pairs is sufficient to ensure local confluence. In our case, this cannot be reduced to checking such minimal states, but the construction is useful as a way to group the cases that need to be considered. We adapt the definition of [5] as follows.

Definition 12. An α₁-critical pattern (with evaluated guards) is of the form

\[
\langle S_1\sigma_1, \{R_1\} \rangle \xrightarrow{\ell_1} \langle S, \emptyset \rangle \xrightarrow{\ell_2} \langle S_2\sigma_2, \{R_2\} \rangle
\]

whenever there exist, for \( k = 1, 2 \) indexed rules \( r_k = (A_k \setminus B_k <--> g_k \mid C_k) \), and application record \( R_k = (r_{\ell_1}i_1^k, \ldots, i_n^k) \) where \( i_1^k, \ldots, i_n^k \) is the list of indices in \( A_k, B_k, \) and \( C_k \), and \( S, S_1 \) and \( S_2 \) are determined in the following way.

- Let \( H_k = A_k \cup B_k, k = 1, 2 \), and split \( B_1 \) and \( H_2 \) into disjoint subsets by \( B_1' = B_1'' \cup B_1''' \) and \( H_2' = H_2'' \cup H_2''' \), where \( B_1'' \) and \( H_2'' \) must have the same number of elements \( \geq 1 \).
- The set of indices used in \( B_1'' \) and \( H_2'' \) are assumed to be identical, and any other index unique, and \( \sigma \) is a most general unifier of \( B_1'' \) and a permutation of \( H_2'' \).
An $\alpha_1$-critical pattern (with delayed guards) is of the form

$$\langle S_1, \{R_1\} \rangle \xrightarrow{\tau_1} \langle S, \emptyset \rangle \xrightarrow{\tau_2} \langle S_2, \{R_2\} \rangle,$$

where all parts are defined as above, except in the last step, that one of $g_k$ contains either an extra-logical built-in or its evaluation leads to a logical built-in $b$ with $\text{Exe}(b) = \text{error}$; the guards $g_k \sigma$ are recognized as the unevaluated guards.

The constraints needed to produce the derivation steps that ensure joinability may not appear in the patterns of Definition 12 but are implied by the invariant. To cope with this, we proceed as follows.

**Definition 13.** An $\alpha_1$-critical corner $y \xrightarrow{\tau_1} x \xrightarrow{\tau_2} y'$ is covered by an $\alpha_1$-critical pattern with evaluated guards of the form $\langle S_1, A_1 \rangle \xrightarrow{\tau_1} \langle S, \emptyset \rangle \xrightarrow{\tau_2} \langle S_2, A_2 \rangle$, whenever there exist a set of indexed constraints $S^+$, a substitution $\sigma^+$ and a set of application records $A^+$ such that $I((S\sigma^+ \cup S^+, A^+))$ holds and

$$y = \langle S_1\sigma^+ \cup S^+, A_1 \cup A^+ \rangle, \quad x = \langle S\sigma^+ \cup S^+, A^+ \rangle, \quad y' = \langle S_2\sigma^+ \cup S^+, A_2 \cup A^+ \rangle.$$

An $\alpha_1$-critical corner $y \xrightarrow{\tau_1} x \xrightarrow{\tau_2} y'$ is covered by an $\alpha_1$-critical pattern with delayed guards of the form $\langle S_1, A_1 \rangle \xrightarrow{\tau_1} \langle S, \emptyset \rangle \xrightarrow{\tau_2} \langle S_2, A_2 \rangle$, whenever there exist a set of indexed constraints $S^+$, a substitution $\sigma^+$ and a set of application records $A^+$ such that $I((S\sigma^+ \cup S^+, A^+))$ holds and

$$y = \langle (S_1\sigma^+ \cup S^+)\text{Exe}(g_1\sigma^+), A_1 \cup A^+ \rangle, \quad x = \langle S\sigma^+ \cup S^+, A^+ \rangle, \quad y' = \langle (S_2\sigma^+ \cup S^+)\text{Exe}(g_2\sigma^+), A_2 \cup A^+ \rangle;$$

$g_1, g_2$ are the unevaluated guards of the pattern.

Analogously to previous results on confluence of CHR, we can state the following.

**Lemma 2.** For a given $I$-terminating program with invariant $I$ and equivalence $\approx$, the set of critical $\alpha_1$-patterns is finite, and any critical $\alpha_1$-corner is covered by some critical $\alpha_1$-pattern.

Our work is currently intended for manual proofs and we do not assume any formal and decidable ways of specifying $I$ and $\approx$, which would be needed for fully mechanizable proof methods. It is straightforward to define critical $\alpha_2$- and $\alpha_3$-patterns and formulate analogous versions of Lemma 2 (left out for reasons of space). These are not needed in our examples as the $\alpha_2$ and $\alpha_3$ properties will appear to be trivial. It is not possible to cover $\beta$-corners by patterns in a similarly syntactic way as two of the participating states are related only semantically by $\approx$. However, we can reuse the developments of and joinability results derived by their methods, e.g., using automatic checkers for classical confluence.
Lemma 3. If a critical $\alpha_1$-pattern $\pi$ (viewed as an $\alpha_1$-corner) is joinable under $I$ and the identity equivalence, any $\alpha_1$-corner covered by $\pi$ is joinable under $I$ and $\approx$.

Example 3 (example 1, ct’d). We formulate invariant and equivalence and prove confluence. The propagation history can be ignored as there are no propagations.

$I$: $I(x)$ holds if and only if $x = \{\text{set}(L)\} \cup \text{Items}$, where $\text{Items}$ is a set of item/1 constraints whose argument is a constant and $L$ a list of constants.

$\approx$: $x \approx x'$ if and only if $x = \{\text{set}(L)\} \cup \text{Items}$ and $x' = \{\text{set}(L')\} \cup \text{Items}$ where $\text{Items}$ is a set of item/1 constraints and $L$ is a permutation of $L'$.

There are no built-ins and thus no critical $\alpha_2$- or $\alpha_3$-patterns. There is only one critical $\alpha_1$-pattern, namely

$\{\text{set}([B|L]), \text{item}(A)\} \leftarrow \{\text{set}(L), \text{item}(A), \text{item}(B)\} \mapsto \{\text{set}([A|L]), \text{item}(B)\}$,

where $L$, $A$, and $B$ are terms such that $I$ holds for the indicate states. Joinability for any corner covered by this pattern is shown by applying the rule to the two “wing” states to form two states $\{\text{set}([B,A,|L])\} \approx \{\text{set}([A,B,|L])\}$.

To check the $\beta$ property, we notice that any $\beta$-corner is of the form

$\{\text{set}(L'), \text{item}(A)\} \cup \text{Items} \approx \{\text{set}(L), \text{item}(A)\} \cup \text{Items} \mapsto \{\text{set}([A|L])\} \cup \text{Items}$

where $L$ and $L'$ are lists, one being a permutation of the other. Applying the rule to the “left wing” state leads to $\{\text{set}([A|L'])\} \cup \text{Items}$ which is equivalent (wrt. $\approx$) to the “right wing” state. As the program is clearly $I$-terminating, it follows that it is confluent modulo $\approx$.

6 Confluence of Viterbi Modulo Equivalence

Dynamic programming algorithms produce solutions to a problem by generating solutions to growing sub-problems, extending those solutions already found. The Viterbi algorithm [11] finds a most probable path of state transitions in a Hidden Markov Model (HMM) that produces a given emission sequence $Ls$, also called the decoding of $Ls$; see [12] for a background on HMMs. There may be exponentially many paths but an early pruning strategy ensures linear time. The algorithm has been studied in CHR by [13], starting from the following program.

:- chr_constraint path/4, trans/3, emit/3.

expand @ trans(Q,Q1,PT), emit(Q,L,PE), path([L|Ls],Q,P,PathRev) ==>
P1 is P*PT*PE  |  path(Ls,Q1,P1,[Q1|PathRev]).

prune @ path(Ls,Q,P1,_) \ path(Ls,Q,P2,_) <=< P1 >= P2 | true.

The meaning of a constraint path($Ls,q,p,R$) is that $Ls$ is a remaining emission sequence to be processed, $q$ the current state of the HMM, and $p$ the probability of a path $R$ found for the already processed prefix of the emission sequence.
To simplify the program, a path is represented in reverse order. Constraint \( \text{trans}(q,q',pt) \) indicates a transition from state \( q \) to \( q' \) with probability \( pt \), and \( \text{emit}(q,\ell,pc) \) a probability \( pc \) for emitting letter \( \ell \) in state \( q \).

Decoding of a sequence \( Ls \) is stated by the query “HMM, \( \text{path}(Ls,q0,1,[]) \)” where HMM is an encoding of a particular HMM in terms of \text{trans} and \text{emit} constraints. Assuming HMM and \( Ls \) be fixed, the state invariant \( I \) is given as reachability from the indicated query. The program is \( I \)-terminating, as a new \text{path} constraint introduced by the expand rule has a first argument shorter than that of its predecessor. Depending on the application order, it may run in between linear and exponential time, and [13] proceeds by semantics preserving program transformations that lead to an optimal execution order.

The program is not confluent in the classical sense, i.e., without an equivalence, as the \text{prune} rule may need to select one out of two different and equally probable paths. A suitable state equivalence may be defined as follows.

**Definition 14.** Let \( (\text{HMM} \cup \text{PATHS}_1,T) \approx (\text{HMM} \cup \text{PATHS}_2,T) \) whenever: For any indexed constraint \((i:\text{path}(Ls,q,P,R_i)) \in \text{PATHS}_1\) there is another \((i:\text{path}(Ls,q,P,R_2)) \in \text{PATHS}_2\) and vice versa.

The built-ins used in guards, \( \text{is}/2 \) and \( \geq/2 \), are logical and complete with respect to \( I \) so there are no \( \alpha_2 \)- or \( \alpha_3 \)-critical corners. For simplicity of notation, we ignore the propagation histories. There are three critical \( \alpha_1 \) patterns to consider:

(i) \( y \leftarrow x \rightarrow y' \), where \( x \) contains two \text{path} constraints that may differ only in their last arguments, and \( y \) and \( y' \) differ only in which of these constraints that are preserved; thus \( y \approx y' \).

(ii) \( y \leftarrow x \rightarrow y' \), where \( x = \{\pi_1,\pi_2,\tau,\eta\}, \pi_i = \text{path}(L,q_i,P_i,R_i) \) for \( i = 1,2, P_1 \geq P_2, \) and \( \tau,\eta \) the \text{trans} and \text{emit} constraints used for the expansion step. Thus \( y = \{\pi_1,\tau,\eta\} \) and \( y' = \{\pi_1,\pi_2,\pi_2',\tau,\eta\} \) where \( \pi_2' \) is expanded from \( \pi_2 \). To show joinability, we show the stronger property of the existence of a state \( z \) with \( y \rightarrow z \leftarrow y' \). We select \( z = \{\pi_1,\pi_1',\tau,\eta\} \), where \( \pi_1' \) is expanded from \( \pi_1 \). The probability in \( \pi_1' \) is greater or equal to that of \( \pi_2' \), which means that a pruning of \( \pi_2 \) is possible when both are present. Joinability is shown as follows.

\[
\begin{align*}
y & \rightarrow z & \leftarrow z & \rightarrow y' \\
\{\pi_1,\pi_1',\pi_2,\tau,\eta\} & \leftarrow \{\pi_1,\pi_1',\pi_2,\tau,\eta\} & \leftarrow \{\pi_1,\pi_1',\pi_2',\tau,\eta\} & \leftarrow \{\pi_1,\pi_1',\pi_2',\tau,\eta\}
\end{align*}
\]

(iii) As case (ii) but with \( P_2 \geq P_1 \) and \( y = \{\pi_2,\tau,\eta\} \); proof similar and omitted.

Thus all \( \alpha \)-critical corners are joinable. There are no critical \( \beta \) corners, as whenever \( x' \approx x \rightarrow y \), the rule \( r \) can apply to \( x' \) with an analogous result, i.e., there exists a state \( y' \) such that \( x' \rightarrow y' \approx y \). This finishes the proof of confluence modulo \( \approx \).

\(^1\) It may be the case that \( \pi_1' \) was produced and pruned at an earlier stage, so the propagation history prevents the creation of \( \pi_1' \) anew. A detailed argument can show, that in this case, there will be another constraints \( \pi_1' \) in the store similar to \( \pi_1' \) but with a \( \geq \) probability, and \( \pi_1' \) can be used for pruning \( \pi_2' \) and obtain the desired result in that way.
7 Confluence of Union-Find Modulo Equivalence

The Union-Find algorithm [14] maintains a collection of disjoint sets under union, with each set represented as a tree. It has been implemented in CHR by [15] who proved it nonconfluent using critical pairs [5]. We have adapted a version from [6], extending it with a new token constraint to be explained; let $UF_{\text{token}}$ refer to our program and $UF_0$ to the original without token constraints.

union @ token, union(A,B) <=> find(A,X), find(B,Y), link(X,Y).
findNode @ A ~> B \ find(A,X) <=> find(B,X).
findRoot @ root(A) \ find(A,X) <=> A=X.
linkEq @ link(A,A) <=> token.
link @ root(A) \ link(A,B), root(B) <=> B ~> A, token.

The ~> and root constraints, called tree constraints, represent a set of trees. A finite set $T$ of ground tree constraints is consistent whenever: for any constant $a$ in $T$, there is either one and only one root($a$) $\in T$, or $a$ is connected via a unique chain of ~> constraints to some $r$ with root($r$) $\in T$. We define sets($T$) to be the set of sets represented by $T$, formally: the smallest equivalence relation over constants in $T$ that contains the reflexive, the transitive closure of ~>; set($a$,T) refers to the set in sets($T$) containing constant $a$.

Our allowed queries are ground and of the form $T \cup U \cup \{\text{token}\}$, where $T$ is a consistent set of tree constraints, and $U$ is a set of constraints union($a_i$,b), where $a_i$, $b_i$ appear in $T$. The invariant $I$ is defined by reachability from these queries. By induction, we can show the following properties of any $I$-state $S$.

- Either $S = T \cup U \cup \{\text{token}\}$, where $T$ is a consistent set of tree constraints and $U$ a set of constraints union($a_i$,b), $a_i$, $b_i$ appear in $T$.

- $S = T \cup U \cup \{\text{link}(A_1,A_2)\} \cup F_1 \cup F_2$ where $T$, $U$ are as in the previous case, and for $i = 1,2$,
  - if $A_i$ is a constant, $F_i = \emptyset$, otherwise
  - $F_i = \{\text{find}(a_i,A_i)\}$ or $F_i = \{(a_i = A_i)\}$ for some constant $a_i$.

As shown by [15], $UF_0$ is not confluent in the classical sense, which can be related to the following issues.

(i) When the detailed steps of two union operations are intertwined in an unfortunate way, the program may get stuck in a state where it cannot finish the operation as shown in the following derivation.

root(a), root(b), root(c), union(a,b), union(b,c) $\mapsto$
root(a), root(b), root(c), link(a,b), link(b,c) $\mapsto$
b ~> a, root(a), root(c), link(b,c)

(ii) Different execution orders of the union operations may lead to different data structures (representing the same sets). This is shown in the following derivations from a query $q_0 = \{\text{root}(a),\text{root}(b),\text{root}(c),\text{union}(a,b),\text{union}(b,c)\}$.

$q_0 \mapsto$ root(a), root(b), b ~> a, union(b,c) $\mapsto$ root(a), b ~> a, c ~> a
$q_0 \mapsto$ root(a), root(b), c ~> b, union(a,b) $\mapsto$ root(b), b ~> a, c ~> b

We proceed, now, to show that $UF_{\text{token}}$ is confluent modulo an equivalence $\approx$, defined as follows.
that rule with itself. We reuse the analysis of [15] who identified all critical pairs for is locally confluent modulo "constraint, it can be seen that there are no critical \( \beta \) inspection of the possible derivation steps one by one (for each rule and for the constraints and formalized in our invariant \( I \) follows that \( UF \) token of implemented CHR systems (which may seem a bit unusual in a context concerned with confluence). In [6], those pairs are avoided by restricting allowed arguments can that, but avoid the problem due to the control patterns imposed by the token constraint. As noticed by [6], these arguments can never occur for a union \( \alpha \) token \( UF \) token reach a state \( T \cup \{ \text{token} \} \approx T' \cup \{ \text{token} \} \) whenever \( sets(T) = sets(T') \).

- \( T \cup U \cup \{ \text{token} \} \approx T' \cup U \cup \{ \text{token} \} \) whenever \( sets(T) = sets(T') \).

- \( T \cup U \cup \{ \text{link}(A_1, A_2) \} \cup F_1 \cup F_2 \approx T' \cup U \cup \{ \text{link}(A'_1, A'_2) \} \cup F'_1 \cup F'_2 \) whenever \( sets(T) = sets(T') \) and for \( i = 1, 2 \), that
  - if \( A_i \) is a constant and (by \( I \)) \( F_i = \emptyset \), then \( A'_i \) is a constant, \( set(A_i, T) = set(A'_i, T') \) and \( F'_i = \emptyset \)
  - if \( A_i \) is a variable and \( F_i = \{ \text{find}(a_i, A_i) \} \) for some constant \( a_i \), then \( F'_i = \{ \text{find}(a'_i, A'_i) \} \) and \( set(a_i, T) = set(a'_i, T') \),
  - if \( A_i \) is a variable, \( F_i = \{ \{ a_i = A_i \} \} \) for some constant \( a_i \) with \( \text{root}(a_i) \in T \) then \( F'_i = \{ \{ a'_i = A'_i \} \} \), \( \text{root}(a'_i) \in T' \) and \( set(a_i, T) = set(a'_i, T') \).

There are no critical \( \alpha \) token and \( \alpha \) token patterns. The \( \alpha \) patterns (critical pairs) of \( UF \) token are those of \( UF_0 \) and a new one, formed by an overlap of the union rule with itself. We reuse the analysis of [15] who identified all critical pairs for \( UF_0 \); by Lemma[3] we consider only those pairs, they identified as non-joinable.

They identified eight non-joinable critical pairs; the first one (“the unavoidable” pair) concerns issue (ii). Its ancestor state \( \{ \text{find}(B, A), \text{root}(B), \text{root}(C), \text{link}(C, B) \} \), is excluded by \( I \): any corner covered, \( B \) and \( C \) must be ground, thus also the link constraint, which according to \( I \) excludes a find constraint. This can be traced to the effect of our token constraint, that forces any union to complete its detailed steps, before a next union may be entered. However, issue (ii) pops up in the new pattern for \( UF \) token, \( y \leftarrow x \mapsto y' \) where:

\[
\begin{align*}
x &= \{ \text{token}, \text{union}(A, B), \text{union}(A', B') \} \\
y &= \{ \text{find}(A, X), \text{find}(B, Y), \text{link}(X, Y), \text{union}(A', B') \} \\
y' &= \{ \text{find}(A', X'), \text{link}(X', Y'), \text{union}(A, B) \}
\end{align*}
\]

To show joinability of any corner covered by this pattern means to find \( z, z' \) such that \( y \mapsto z \approx z' \mapsto y' \). This can be done by, from \( y \), first executing all remaining steps related to \( \text{union}(A, B) \) and then the steps relating to \( \text{union}(A', B') \) to reach a state \( z = T \cup U \cup \{ \text{token} \} \). In a similar way, we construct \( z' = T' \cup U \cup \{ \text{token} \} \), starting with the steps relating to \( \text{union}(A', B') \) followed by those of \( \text{union}(A, B) \). It can be proved by induction that \( sets(T) = sets(T') \), thus \( z \approx z' \).

Next, [15] identifies three critical pairs, that imply inconsistent tree constraints. The authors argue informally that these pairs will never occur for a query with consistent tree constraints. As noticed by [6], these arguments can be formalized using an invariant. The last four pairs of [15] relate to issue (i) above; [15] argues these to be avoidable when assuming procedural properties of implemented CHR systems (which may seem a bit unusual in a context concerned with confluence). In [6], those pairs are avoided by restricting allowed queries to include only a single union constraint; we can allow any number of those, but avoid the problem due to the control patterns imposed by the token constraints and formalized in our invariant \( I \).

This finishes the argument that \( UF \) token satisfies the \( \alpha \) property, and by inspection of the possible derivation steps one by one (for each rule and for the "=" constraint), it can be seen that there are no critical \( \beta \) corners. Thus \( UF \) token is locally confluent modulo \( \approx \), and since tree consistency implies termination, it follows that \( UF \) token is confluent modulo \( \approx \).
8 Discussion and detailed comments on related work

We already commented on the foundational work on confluence for CHR by [5], who, by the use of Newman’s lemma, could devise a method to prove confluence by inspecting a finite number of critical pairs. This formed also the foundation of automatic confluence checkers [5,7,10] (with no invariant and no equivalence).

The addition of an invariant $I$ in the specification of confluence problems for CHR was suggested by [6]. The authors considered a construction similar to our $\alpha_1$-corners and critical $\alpha_1$-patterns. They noted that critical $\alpha_1$-patterns usually do not satisfy the invariant, so they based their approach on defining a collection of corners based on $I$-states as minimal extensions of such patterns. Local confluence, then, follows from joinability of this collection of minimally extended states. However, there are often infinitely many such minimally extended states; this happens even for a natural invariant such as groundness when infinitely many terms are possible, as is the case in Prolog based CHR versions. We can use this construction (in cases where it is finite!) to further cluster the space of our critical corners, but our examples worked quite well without this.

Of other work concerned with confluence for CHR, we may mention [16] who considered confluence for non-terminating CHR programs, used recently by [17] for a specific type inference problem. We may also refer to [18] that gives an overview of CHR related research until 2010, including on confluence.

9 Conclusion and future work

We have introduced confluence modulo equivalence for CHR, which allows a much larger class of programs to be characterized as confluent in a natural way, thus increasing the practical relevance of confluence for CHR.

We demonstrated the power of the framework by showing confluence modulo equivalence for programs that use a redundant data representation (the set-as-lists and Union-Find programs) and a dynamic programming algorithm (the Viterbi program); all these are out of scope of previous confluence notions for CHR. With the new operational semantics, we can also handle extra-logical and incomplete built-in predicates, and the notational improvements obtained by this semantics may also promote new research on confluence.

As a first steps towards semi- or fully automatic proof methods, it is important to notice that classical joinability of a critical pair – as can be decided by existing confluence checkers such as [10] – provide a sufficient condition for joinability modulo any equivalence.

Thus only classically non-joinable pairs – in our terminology $\alpha_1$ patterns – need to be examined in more details involving the relevant equivalence; however, in some cases there may also be critical $\alpha_2$, $\alpha_3$ and $\beta$ patterns that need to be considered. Machine supported proof methods for these parts need formal and decidable specifications of (i) invariant and (ii) equivalence, which may be supplied by additional CHR programs. For (i), it may be done by checking for violations, and for (ii), such a program may normalize states into a canonical form that can be compared in a straightforward way.
References

Exhaustive Execution of CHR through Source-to-Source Transformation

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Abstract. Constraint Handling Rules (CHR) is a committed-choice rule-based programming language. Rules rewrite a global multi-set of constraints to another. Overlapping sets of constraints within the rules and the order of constraints within rules and queries entail different derivation paths. In this work, a novel operational strategy is proposed which enables a high-level form of execution control that empowers a comprehensive and customizable execution strategy. It allows full space exploration for any CHR program, thus finding all possible results to a query which is interesting for many non-confluent programs. The proposed transformation is performed as a source-to-source transformation from any CHR program to one utilizing disjunction to force an exhaustive explorative execution strategy. The work is complemented by formal arguments to prove the correctness and completeness of the transformation.

Keywords: Constraint Handling Rules, Execution flow control, Exhaustive execution, Search space exploration, Source-to-source transformation

1 Introduction

Constraint Handling Rules (CHR) is a rule-based programming language based on a set of multi-headed guarded rewrite rules \cite{9}. The rules operate on a global multi-set of constraints to rewrite them from one multi-set of constraints to another. The rules are applied exhaustively until a final state is attained; where no more rules are applicable. The execution of a program on an initial query is known as a derivation. Derivations follow a committed-choice manner, where a chosen rule cannot be retracted. Rules could operate on overlapping sets of constraints. The order of constraints within the rules, the order of constraints within the input query, and the actual order of program rules chosen for execution are factors defined by the operational semantics that determine the derivation path.

For a given CHR program and for any initial state, if all derivations from that state result in the same final state then the program is confluent. On the other hand, non-confluent programs produce different final states for a given query depending on the chosen derivation path. Non-confluent programs are common specially with agent-oriented programming. The Blocks World example \cite{8, 12}
can be modelled in CHR to describe the behaviour of an agent in a world of objects.

**Example 1. Blocks World**

```
rule-1 @ get(X), empty <=> hold(X).
rule-2 @ get(X), hold(Y) <=> hold(X), clear(Y).
```

The constraint `hold(X)` represents that the agent holds an object `X`. The constraint `empty` represents that the agent holds nothing, the constraint `clear(Y)` denotes that the object `Y` is not held, and the constraint `get(X)` represents the agent’s action to get an object `X`. The program consists of two simplification rules having overlapping heads, both rules have a `get(X)` constraint on the left-hand side. The simplification rules rewrite matched left-hand side constraints to the right-hand side constraints. With the abstract operational semantics and an initial query ‘`empty, get(box), get(cup)`’, two disjoint final states can be reached:

![Derivation tree for Example 1 with query: empty,get(box),get(cup).](image)

The CHR implementation of the K.U. Leuven system [10], which follows the refined operational semantics but with additional fixed orders, entails that only one final state is reached ‘`hold(cup), clear(box)`’. However, changing the order of program rules, the order of constraints within an input query or within the rules, yields different final states for a non-confluent program. For example, executing the program with the reordered query: ‘`empty, get(cup), get(box)`’ entails the final state ‘`hold(box), clear(cup)`’.

In order to change the program execution flow, source-to-source transformations are used to facilitate a straightforward implementation on top of existing CHR implementations whilst exploiting the optimizations of current CHR compilers [14]. For example, such transformations can be used to transform CHR with user-defined priorities into regular CHR [4, 6] and to a tree-based semantics with a breadth-first search [5].

In this work, we capture the essence of the tree-based semantics [5] by introducing branching derivations through identifying overlapping rule-heads. The tree-based semantics makes use of the disjunctive branching points within a program to generate CHR constraints to encode nodes and edges.
The presented approach is based on transforming the input program to a program that utilizes Constraint Handling Rules with Disjunction (CHR $\lor$); which is an extension of CHR [3]. CHR $\lor$ facilitates relaxing the committed-choice behaviour of CHR by allowing disjunctive bodies and hence backtracking over alternatives. Application of a transition rule generates a branching derivation which can be utilized to overcome the committed-choice execution of CHR.

This work introduces an exhaustive execution strategy as a source-to-source transformation of the CHR program rules. It finds all possible results to a query for non-confluent programs. This execution platform is extensible to adopt different search strategies on a program’s search space. For implementing constraint solvers, this could definitely enable optimized search techniques [5]. The angelic semantics of CHR [13] is a preliminary similar work aimed to explore all possible execution choices using derivation nets. However, no implementation nor definition of an operational semantics was provided. The transformation proposed here aims to reach an angelic semantics however by generating search trees similar to the tree-based semantics without explicitly encoding edges and nodes.

In this work, we establish a strong correspondence of the proposed transformation to the operational semantics of CHR to prove the soundness and completeness of the proposed transformation.

Moreover, as a practical application of this work, the exhaustive execution by transformation is used as a top-level execution strategy required for the inverse execution of CHR [15]. The inverse execution of CHR presents a simple inversion transformation, that changes a CHR program to its inverse one by interchanging rule sides. For the blocks world example mentioned earlier, that would mean a program which given an output state ‘hold(cup), clear(box)’ generates all possible input states that lead to it. With a normal CHR execution implementation only one input is uncovered; thus a different exhaustive execution strategy is required. A primitive unformalized execution transformation was proposed in [15]. In this work, we extend the work of [15] by using the proposed transformation for executing inverse CHR programs.

The paper proceeds by providing some background information in Section 2, followed by the proposed transformation in Section 3. The proof of the soundness and completeness of the transformation is given in Section 4. Then an application of the proposed transformation is given in Section 5, followed by some concluding remarks in Section 6.

## 2 Constraint Handling Rules

### 2.1 Syntax

Constraint Handling Rules (CHR) [9, 10] consist of guarded rewrite rules that perform conditional transformation of multi-sets of constraints, known as a constraint store. There are three types of rules: simplification, propagation and simpagation rules. All rules have an optional unique identifier $\text{rule-id}$ separated from the rule body by @. For all rules, $H_k$ and/or $H_r$ form the rule head. $H_k$ and
$H_r$ are sequences of user-defined CHR constraints that are kept and removed from the constraint store respectively. $Guard$ is the rule guard consisting of a sequence of built-in constraints, and $B$ is the rule body comprising of CHR and built-in constraints. The generalized simpagation rules are of the form:

$$rule-id @ H_k \setminus H_r \iff Guard \mid B$$

The two other types of rules are special cases of the generalized simpagation rule. Simplification rules have no kept head constraints $H_k$ while propagation rules do not remove any head constraints. These rules are of the form:

Simplification rule: $rule-id @ H_r \iff Guard \mid B$  
Propagation rule: $rule-id @ H_k \implies Guard \mid B$

Propagation and simpagation rules can be represented as simplification rules by adding kept constraints to the rule bodies; in its abstract sense this introduces the notion of non-termination. However this issue can be avoided by using a refined operational semantics which uses a token store as a history for propagated constraints [1]. For simplicity, in this paper the focus will be only with simplification rules, however the work can be extended in a straightforward manner by adding the notion of a propagation history.

Constraint Handling Rules with Disjunction (CHR$^\lor$) [3] is an extension which allows disjunctive rule bodies and enables a backtrack search over alternatives. A simplification rule with two possible bodies ($B_1 \lor B_2$) is of the form:

$$ruleV-id @ H_r \iff Guard \mid B_1 \lor B_2$$

### 2.2 Operational Semantics

An operational semantics of CHR describes the behaviour of a program in terms of a state transition system which models the execution of a program.

The abstract operational semantics $\omega_{va}$ of CHR is formulated as a state transition system, where states are goals comprising a conjunction of built-in and CHR constraints, and a transition defines an execution state and its subsequent execution state. An initial state is an arbitrary one and a final state is a terminal one where no further transitions are possible. The abstract operational semantics includes one transition rule shown below, where $P$ is a CHR program with rules $r$ and $CT$ is the constraint theory for the built-in constraints.

**Apply**

$$\begin{align*}
(H_k \land H_r \land C) & \rightarrow^r_{apply} (H_k \land Guard \land B \land C) \\
\text{if there is an instance of a rule } r \text{ in } P \text{ with new local variables } \bar{x} \text{ such that:} \\
& r @ H_k \setminus H_r \iff Guard \mid B \text{ and } CT \models \forall(C \rightarrow \exists \bar{x} \ Guard)
\end{align*}$$

The extended operational semantics for CHR$^\lor$ operates on a disjunction of CHR states known as a configuration: $s_1 \lor s_2 \lor \cdots \lor s_n$. The semantics includes two transitions; one containing the original apply transition which is applicable
to a single state. The second transition is a split transition which is applicable to any goal containing a disjunction. It leads to a branching derivation entailing two states, where each state can be processed independently. The second split transition required by CHR $\lor$ is shown below.

\[
\text{Split} \quad ((H_1 \lor H_2) \land C) \lor S \rightarrow_{\text{split}} (H_1 \land C) \lor (H_2 \land C) \lor S
\]

Such a derivation can be visualized as a tree; consisting of a set of nodes and a set of directed edges connecting these nodes [11]. The root node corresponds to the initial state or goal. An intermediate node represents a midway non-final state within a derivation. A leaf node represents a successful or failed final state. Edges between nodes denote the applied transitions. Due to the clarity offered by search trees for such derivations, in this work all examples will be depicted using trees.

The refined operational semantics $\omega_r$ is a theoretical operational semantics that makes execution considerably more deterministic by establishing an order for goal constraints [7]. It is more closely related to actual implementations of CHR. Due to the limited space of this paper, we refer to the $\omega_r$ as defined in [7].

3 Transformation

In this section, we describe how to transform a CHR program $P$ to a CHR program $P^T$, such that execution of program $P^T$ under the refined operational semantics $\omega_r$ produces all possible final results. The transformation is divided into three sets of rules:

**Definition 1. Modified Source Rules** For every rule $r@H \Leftrightarrow \text{Guard} \mid B$ in program $P$ with $H = c_1(X_{11},...,X_{1n_1}),...,c_l(X_{l1},...,X_{ln_l})$, a modified source rule $r_t$ is added to the transformed program, defined as follows:

\[
\begin{align*}
\text{Modified Source Rules} \\
r_t & \equiv \text{depth}(Z), c_1^T(X_{11},...,X_{1n_1},y_1,Z),...,c_l^T(X_{l1},...,X_{ln_l},y_l,Z) \\
& \Leftrightarrow \text{Guard} \mid B, \text{depth}(Z + 1)
\end{align*}
\]

As a result, every modified source rule $r_t$ in $P^T$ corresponds to rule $r$ in $P$. Every constraint $c(X_1,...,X_n)$ in the head of rule $r$ is transformed to constraint $c^T(X_1,...,X_n,y,Z)$ in rule $r_t$. The argument $y$ represents the $y$th occurrence of constraint $c/n$. Let $m$ be the number of occurrences of $c/n$ in the source program, in the transformed program the first occurrence of $c(X_1,...,X_n)$ is transformed to $c^T(X_1,...,X_n,1,Z)$, the second occurrence of $c(X_1,...,X_n)$ is transformed to $c^T(X_1,...,X_n,2,Z)$, and the last occurrence of $c(X_1,...,X_n)$ is transformed to $c^T(X_1,...,X_n,m,Z)$. The argument $Z$ and constraint $\text{depth}/1$ are explained in the Assignment Rules.

**Definition 2. Assignment Rules** For every constraint $c(X_1,...,X_n)$ that appears in a rule head in program $P$, an assignment rule is added to the transformed program, defined as follows:
Assignment Rules

\[ \text{assign}_c \oplus \text{depth}(Z) \setminus c(X_1, \ldots, X_n) \]
\[ \leftrightarrow c^T(X_1, \ldots, X_n, 0, Z) \lor \ldots \lor c^T(X_1, \ldots, X_n, m, Z) \]

An assignment rule simplifies constraint \( c(X_1, \ldots, X_n) \) to constraint \( c^T(X_1, \ldots, X_n, Y, Z) \). The argument \( Y \) is used to have at most one possible matching between the transformed constraint and the constraints that appear in the heads of the Modified Source Rules. The domain of \( Y \) is from 0 to \( m \), in which \( m \) is the number of occurrences of constraint \( c/n \) in the rule heads of program \( P \). The argument 0 is used to prevent an active constraint \( c^T(X_1, \ldots, X_n, 0, Z) \) from being matched during its current state in the tree, thus allowing it to be matched later in the derivation tree. Constraint \( \text{depth}/1 \) represents the depth of a state in the derivation tree. The depth is increased when a modified source rule is fired.

Definition 3. Reset Rules For every constraint \( c(X_1, \ldots, X_n) \) that appears in a rule head in program \( P \), a reset rule is added to the transformed program, defined as follows:

\[ \text{reset}_c \oplus \text{depth}(Z) \setminus c^T(X_1, \ldots, X_n, 0, Z') \]
\[ \leftrightarrow Z' < Z \lor c(X_1, \ldots, X_n) \]

The purpose of reset rules is to reset the unmatched constraints \( c^T(X_1, \ldots, X_n, 0, Z') \) if a newly state in the tree is derived. The constraint \( c^T(X_1, \ldots, X_n, 0, Z') \) will be reset to constraint \( c(X_1, \ldots, X_n) \), thus allowing it to be re-assigned by the Assignment Rules.

Finally, the rule \( \text{start} \leftrightarrow \text{depth}(0) \) is added to the transformed program to trigger the constraints needed for execution. Moreover, a constraint \( \text{start} \) is added at the end of the initial query.

Example 1 (continued). The transformation of program \( P \):

\[ \text{start} \leftrightarrow \text{depth}(0), \]

Reset Rules

\[ \text{reset}_c \oplus \text{depth}(Z) \setminus c^T(X_1, \ldots, X_n, 0, Z') \]
\[ \leftrightarrow Z' < Z \lor c(X_1, \ldots, X_n) \]

Assignment Rules

Assign empty @ depth(Z) \& empty_t(0,Z1) \leftrightarrow Z1 < Z | empty.
Assign get @ depth(Z) \& get_t(X,0,Z1) \leftrightarrow Z1 < Z | get(X).
Assign hold @ depth(Z) \& hold_t(X,0,Z1) \leftrightarrow Z1 < Z | hold(X).

Modified Source Rules

\[ \text{rule-1}_t \oplus \text{depth}(Z), \text{get}_t(X,1,Z), \text{empty}_t(1,Z) \]
\[ \leftrightarrow \text{hold}(X), Z1 \text{ is } Z+1, \text{depth}(Z1). \]
\[ \text{rule-2}_t \oplus \text{depth}(Z), \text{get}_t(X,2,Z), \text{hold}_t(Y,1,Z) \]
\[ \leftrightarrow \text{hold}(X), \text{clear}(Y), Z1 \text{ is } Z + 1, \text{depth}(Z1). \]
Figure 2 shows the derivation tree of the transformed program when executed by the query \(\text{‘empty, get(box), get(cup), start’}\). The redundant states in the tree are represented by grey arrows. States are duplicated in the tree by the Assignment Rules, because these rules create variants to the same state. For example, the two states \(\text{‘depth(0), empty, get(box,0), get(cup,0)’}\) and \(\text{‘depth(0), empty, get(box,0), get(cup,0)’}\) are variants of the same state.

![Derivation tree for Example 1 after transformation.](image)

Every state in the derivation tree of the source program in Figure 1, has a corresponding result (final state) in the transformed program. Thus, we extend the transformation with pruning rules to remove every result in the transformed program that corresponds to an intermediate (non-final) state in the source program. To map a state in the transformed program to a state in the source program, every constraint \(c^T(X_1, ..., X_n, \_\_\_)\) is mapped to constraint \(c(X_1, ..., X_n)\) and constraint \(\text{depth/1}\) is removed. This mapping can be done by adding bridge rules [2] to the transformed program to relate constraints between the source and transformed programs. These rules are not described here due to their clarity.

**Definition 4. Intermediate States in P**

Let \(S\) be a state in program \(P\), \(S\) is intermediate if there exists a derivation \(S \rightarrow S'\).

**Definition 5. Pruning Rules**

A result \(S^T\) in program \(P^T\) has to be removed if it corresponds to an intermediate state in \(P\). To remove \(S^T\), for every rule \(r \in H \Leftrightarrow \text{Guard} \mid B\) in the source program \(P\), a corresponding pruning-\(r\) rule is added to the transformed program \(P^T\), defined as follows:

\[
\text{Pruning Rules} \\
pruning-r \in H^T \Leftrightarrow \text{Guard} \mid \text{fail}
\]

Every constraint \(c(X_1, ..., X_n)\) in \(H\) is transformed to constraint \(c^T(X_1, ..., X_n, \_\_\_)\) in \(H^T\). Moreover, the rule \(\text{start} \Leftrightarrow \text{depth(0)}\) is transformed to \(\text{start} \Leftrightarrow \text{depth(0), end}\), therefore constraint \(\text{end/0}\) will become active only when all Reset Rules, Assignment Rules, and Modified Source Rules cannot be triggered anymore. If a Pruning Rule can be fired on state \(S^T\) and the \(\text{Guard}\) holds, this implies that rule \(r\) can be applied on state \(S\) that corresponds to state \(S^T\); state \(S^T\) is removed by having a derivation \(S^T \rightarrow \text{fail}\).
Example 1 (transformation with pruning rules).

**start** <=> depth(0) will be modified to **start** <=> depth(0), **end**.

Moreover the following rules will be added to the transformed program:

- **pruning-rule-1** @ **end**, get_t(X,_,_), empty_t(_,_) <=> fail.
- **pruning-rule-2** @ **end**, get_t(X,_,_), hold_t(Y,_,_) <=> fail.

The derivation tree for the query ‘empty, get(box), get(cup), start’ is depicted in Figure 3. As shown in the figure, every result that corresponds to an intermediate state in Figure 1 fails after adding the pruning rules to the transformed program.

**Fig. 3.** Derivation tree for Example after pruning.

4 Soundness and Completeness of the Transformation

In this section, soundness and completeness of the transformation are proved. The abstract operational semantics will be used for the derivations throughout the proofs. However for the completeness proof, derivations of the transformed program are based on the refined operational semantics \( \omega_r \), because the current CHR compiler is based on this semantics. Therefore it proves that the transformed program produces all possible results when executed under the semantics of the CHR compiler.

**Definition 6.** *Equivalence of two CHR states.* Let \( S \) be a CHR state derived from the initial query \( G \rightarrow^* S \) in program \( P \), and \( W = S_1 \lor \ldots \lor S_n \) be a configuration derived from the initial query \( (G \land \text{depth}(0)) \rightarrow^* W \) in program \( P^T \). CHR state \( S \) is equivalent to CHR state \( S' \) in \( \{S_1, \ldots, S_n\} \) according to the following inductive definition. \( S \equiv S' \) if

**base case:**
- \((S = G_{\text{builtin}}) \land (S' = (G_{\text{builtin}} \land \text{depth}(Z)))\), in which \( G_{\text{builtin}} \) is a conjunction of built-in constraints.

**inductive step:**
- \((S = (c(X_1, \ldots, X_n) \land S_r)) \land (S' = (c^T(X_1, \ldots, X_n, Y, Z) \land S'_r)) \land (S_r \equiv S'_r)\)
- \((S = (c(X_1, \ldots, X_n) \land S_r)) \land (S' = (c(X_1, \ldots, X_n) \land S'_r)) \land (S_r \equiv S'_r)\)
Definition 7. **Soundness of a CHR configuration.** Given an initial query \( G \) and a configuration \( W = S_1 \lor ... \lor S_n \) derived from \( (G \land \text{depth}(0)) \to^* W \) in program \( P^T \). \( W \) is sound if and only if for every CHR state \( S_i \) in \( \{S_1, ..., S_n\} \) there exists a derivation in program \( P \), where \( G \to^* S \) and \( S \equiv S_i \).

**Theorem 1.** Given a CHR Program \( P \), its corresponding transformed program \( P^T \), and an initial query \( G \). \( P^T \) is sound with respect to \( P \), if and only if every derived configuration \( W \) in program \( P^T \) from the query \( (G \land \text{depth}(0)) \) is sound.

**Proof.**

**Base Case.** \( S = G, W = (G \land \text{depth}(0)) \), in which \( G \) is the initial query. Since \( S \equiv W \), then the initial configuration \( W \) is sound.

**Induction Hypothesis.** Let \( W = S_1 \lor ... \lor S_k \) be a configuration derived from \( (G \land \text{depth}(0)) \to^* W \), assume that \( W \) is sound.

**Induction Step.** We prove that if \( W \to W' \), then \( W' \) is sound.

Without loss of generality, we assume that the fired rule is applied on the first CHR state \( S_1 \) in \( W \).

**Case 1. A reset rule is applicable on \( W \) and \( CT \models G_{\text{builtin}} \to (Z' < Z) \)**

\[
W = ((\text{depth}(Z) \land c'(X_1, ..., X_n, 0, Z') \land C) \lor (S_2 \lor ... \lor S_k)) \quad \text{Because a reset rule is applicable on } W
\]

\[
W' = (\text{depth}(Z) \land c(X_1, ..., X_n) \land C \lor (S_2 \lor ... \lor S_k)) \quad \text{Firing a reset rule on } W, \ W \Rightarrow_{\text{reset}} W'
\]

\[
S \equiv (\text{depth}(Z) \land c(X_1, ..., X_n) \land C) \quad \text{According to the hypothesis there exists a derivation } G \to^* S
\]

Therefore \( W' \) is sound.

**Case 2. An assign rule is applicable on \( W \)**

\[
W = ((c(X_1, ..., X_n) \land \text{depth}(Z) \land C) \lor (S_2 \lor ... \lor S_k)) \quad \text{Because an assign rule is applicable on } W
\]

\[
W'' = ((c'(X_1, ..., X_n, 0, Z) \lor ... \lor c'(X_1, ..., X_n, l, Z)) \land \text{depth}(Z) \land C \lor (S_2 \lor ... \lor S_k)) \quad \text{Firing an assign rule on } W, \ W \Rightarrow_{\text{assign}} W''
\]

\[
W''' = ((c'(X_1, ..., X_n, 0, Z) \land \text{depth}(Z) \land C) \lor ... \lor ((c'(X_1, ..., X_n, l, Z) \land \text{depth}(Z) \land C) \lor (S_2 \lor ... \lor S_k)) \quad \text{A split is applied on } W'' \Rightarrow_{\text{split}} W'''
\]

\[
S \equiv (\text{depth}(Z) \land c(X_1, ..., X_n) \land C) \quad \text{According to the hypothesis there exists a derivation } G \to^* S
\]

\[
S = (c(X_1, ..., X_n) \land C') \quad \text{According to the definition of equivalence}
\]

Therefore \( W''' \) is sound.

\[
\forall_{0 \leq i \leq l} ((c'(X_1, ..., X_n, i, Z) \land \text{depth}(Z) \land C) \equiv S)
\]
Case 3. A Modified rule \( r_t \) is applicable on \( W \) and \( CT \models G_{\text{builtin}} \rightarrow Guard \)

\[
W = ((\text{depth}(Z) \land \land c_1(X_{11}, ..., X_{1n_1}, Y_1, Z) \land ... \land c_1(X_{11}, ..., X_{1n_1}, Y_1, Z) \land C) \vee (S_2 \lor \lor S_k))
\]

Because a modified rule \( r_t \) is applicable on \( W \)

\[
S \equiv (\text{depth}(Z) \land \land c_1(X_{11}, ..., X_{1n_1}, Y_1, Z) \land C)
\]

According to the hypothesis there exists a derivation \( G \models S \)

\[
S \equiv (c_1(X_{11}, ..., X_{1n_1}) \land ... \land c_1(X_{i1}, ..., X_{i}) \land C') \land C' \equiv (C \land depth(Z))
\]

Firing rule \( r_t \) on \( W, W \rightarrow^* W' \)

\[
(W' = ((B \land \land depth(Z + 1) \land Guard \land C) \vee (S_2 \lor \lor S_k)) \quad \text{Firing rule} r \quad S \rightarrow^* S'
\]

Therefore \( W' \) is sound

Because \( (B \land \land depth(Z + 1) \land Guard \land C) \equiv S' \)

\( \square \)

For simplicity, we define a state in \( \omega_r \) as the tuple \( (A, S, B) \) [7]. The stack \( A \) is a sequence of built-in constraints and CHR constraints. The store \( S \) is a set of CHR constraints. The built-in store \( B \) is a conjunction of built-in constraints.

**Definition 8.** Let \( S \) be a CHR state derived from the initial query \( G \rightarrow^* S \) in program \( P \) and \( S^T = (A, C, B) \) be a state derived from the initial state \( \langle G++[\text{depth}(0)], \phi, \text{true} \rangle \) in program \( P^T \), in which \( ++ \) is the append operator defined for sequences. Let \( E \) be the property defined on \( S \) and \( S^T \), \( E(S, S^T) \) holds iff \( A = [\text{depth}(Z)] \), \( \text{depth}(Z) \in C \), \( C \land B \equiv S \), and \( Z' < Z \) holds for all transformed constraints \( c^T(X_1, ..., X_n, Y, Z') \) in \( C \).

**Theorem 2.** Given a CHR Program \( P \), its corresponding transformed program \( P^T \), and an initial query \( G \). \( P^T \) is complete if and only if for every CHR state \( S \) derived from the initial query \( G \rightarrow^* S \) in program \( P \), there exists a state \( S^T \) derived from the initial state \( \langle G++[\text{depth}(0)], \phi, \text{true} \rangle \) in program \( P^T \) such that \( \langle G++[\text{depth}(0)], \phi, \text{true} \rangle \rightarrow^* S^T \) and \( E(S, S^T) \) holds.

**Proof.**

**Base Case.** There exists a derivation to state \( S^T \) in \( P^T \) and \( E(G, S^T) \) holds.

\[
\langle G++[\text{depth}(0)], \phi, \text{true} \rangle \rightarrow^* \langle [\text{depth}(0)], G_c, G_b \rangle
\]

All CHR constraints in \( G \) will be activated, added to the CHR store, then dropped from the stack, and all built-in constraints will be solved and removed from the stack.

\[
\rightarrow^* \text{activate depth}(0) \text{ becomes the active constraint}
\]

Since \( E(G, \langle [\text{depth}(0)], \{\text{depth}(0) \} \cup G_c, G_b \rangle) \) holds, then the base case holds.

**Induction Hypothesis.** Let \( S \) be a state in \( P \) derived from the initial query \( G \) and \( S^T \) be a state in \( P^T \) derived from the initial state \( \langle G++, \phi, \text{true} \rangle \), assume that \( E(S, S^T) \) holds.
**Induction Step.** We prove that for a derivation \( S \Rightarrow S' \), there exists a derivation \( S_T \Rightarrow^* S_{1T} \) and \( E(S', S_{1T}) \) holds. For simplicity of the proof, we define the following:

- Let \( r \oplus H \Leftrightarrow \text{Guard} \mid B \) be a rule in program \( P \) that applies on \( S = H \land H_{\text{rest}} \) such that \( S = (H \land H_{\text{rest}}) \mapsto_{\text{apply}} S' = (H_{\text{rest}} \land B \land \text{Guard}) \).
- Let \( \text{map} \) be a function that maps every constraint in \( H \) to a partner constraint in \( S \).
- Let \( r_1 \oplus \text{depth}(Z), H^T \Leftrightarrow \text{Guard} \mid B, \text{depth}(Z + 1) \) be a modified source rule in program \( P^T \), that corresponds to rule \( r \) in program \( P \).
- Let \( \text{occurrence} \) be a function that is applied on a transformed constraint \( c^T \), such that \( \text{occurrence}(c^T(X_1, ..., X_n, Y, Z)) = Y \).
- Let \( \text{dep} \) be a function that is applied on a transformed constraint \( c^T \), such that \( \text{dep}(c^T(X_1, ..., X_n, Y, Z)) = Z \).
- Let \( \text{assign} \) be a function that maps a constraint \( c/n \) in \( S \) to an integer value according to the following:
  - if there exists a constraint \( d \) in \( H \) such that \( \text{map}(d) = c \), then \( \text{assign}(c) = \text{occurrence}(d^T(X_1, ..., X_n, Y, Z)) \), in which \( d^T \) is the transformed constraint in \( H^T \) that corresponds to constraint \( d \)
  - otherwise, \( \text{assign}(c) = 0 \)

\[
S^T = \langle \text{depth}(Z) \rangle \cup H_e \cup H_{\tau, \tau}, H_b \quad \text{According to the hypothesis } \text{dep}(c^T) < Z \text{ for all constraints in } H_{\tau, \tau} \text{ and } (\text{depth}(Z) \land H_e \land H_{\tau, \tau} \land H_b) \equiv S \text{ holds.}
\]

**Reset Rule**

\[
\text{reset} \quad \langle \text{depth}(Z) \rangle \cup H_e \cup \{c\} \cup H_{\tau, \tau} \setminus \{c^T\}, H_b
\]

A \( \text{reset} \) rule will be applied on constraint \( c^T \) in \( H_{\tau, \tau} \), constraint \( c^T \) will be removed from the CHR store and its corresponding constraint \( c \) will be added to the top of the stack, constraint \( c \) will be active and will fire rule \( \text{assign} \). Due to the existence of disjunction in the body of the \( \text{assign} \) rule, a split will be applied. The path that will lead to a correct state derivation is the one with \( \text{occurrence}(c^T) = \text{assign}(c) \). For simplicity the reset of disjuncts will be removed from the stack.

**Assign Rule**

\[
\text{assign} \quad \langle \text{depth}(Z) \rangle \cup H_e \cup H_{\tau, \tau} \setminus \{c^T\}, H_b
\]

\[
\text{assign} \quad \langle \text{depth}(Z) \rangle \cup H_e \cup H_{\tau, \tau} \setminus \{c^T\}, H_b
\]

\[
\text{assign} \quad \langle \text{depth}(Z) \rangle \cup H_e \cup H_{\tau, \tau}, H_b
\]

\[
\text{assign} \quad \langle \text{depth}(Z) \rangle \cup H_e \cup H_{\tau, \tau} \setminus \{c^T\}, H_b
\]

\[
\text{assign} \quad \langle \text{depth}(Z) \rangle \cup H_e \cup H_{\tau, \tau}, H_b
\]

The two rules \( \text{reset} \) and \( \text{assign} \) will be applied repeatedly till all constraints \( c^T \) with \( \text{dep}(c^T) < Z \) in \( H_{\tau, \tau} \) are reset and then assigned, or till rule \( r_1 \) can be fired. The second case is covered by the first one, thus only the first case is considered.
There are two sources of non-determinism in the refined operational semantics in contrast with the actual implementation of CHR. The first source of non-determinism is the order of the constraints added to the top of the stack when the transition Solve+Wake is fired, the completeness proof assumes only that the constraints are added to the top of the stack, therefore, the order of the constraints does not matter. The second source of non-determinism is choosing the partner constraints from the CHR store when the transition Apply is fired, in the completeness proof Assignment Rules and Reset Rules are applied on arbitrary constraints $c$ and $c^T$ respectively. Moreover, the completeness proof shows that if a rule $r$ is applicable on CHR state $S$, then there exists a path in the derivation tree of the transformed program, such that every constraint in the head of rule $r_1$ has only one matching in the CHR store and the rest of the constraints in state $S$ are assigned the occurrence 0, which implies that the only applicable Modified Source Rule is rule $r_t$ and that there is only one way to match the constraints in the CHR store with the constraints in the head of rule $r_t$. Therefore, it is implied that the transformed program produces all final results when executed under the semantics of the CHR compiler, although the proof is based on the refined operational semantics.
5 Inverse Execution using Proposed Transformation

The inverse execution of CHR presented in [15] presents a simple inversion trans-
formation, that changes a forward CHR program to an inverse one by interchang-
ing rule sides. Thus for any forward simplification rule \((H_r \iff \text{Guard} \mid B_b, B_c)\),
where \(B_b\) are the built-in body constraints and \(B_c\) are the CHR body constraints
and \(B_b \cup B_c = B\), an inverse rule would be of the form:

\[
\text{inv-simpf} \in B_c \iff B_b, \text{Guard} \mid H_r
\]

The inverse programs are then transformed using the proposed approach in
order to achieve an exhaustive traversal of the inverse program that reaches all
possible input states. However, since all intermediate states are also possible
input states, thus pruning rules are removed from the transformed program. In
some applications (such as for fault analysis) this step may not be necessary, as
retaining pruning rules results in only uncovering the root output states.

Example 2. Blocks World in Reverse - If one wishes to run the blocks
world (Example 1) in reverse, to retrace an agent’s steps, then we require the
inverse transformation of the program code. The forward run for the query
\(\text{empty, get(box), get(cup)}\) was previously depicted in Figure 1. Using the simple
inversion transformation [15], the inverse program becomes:

\[
\begin{align*}
\text{inv-rule-1} & \iff \text{hold}(X) \Rightarrow \text{get}(X), \text{empty}. \\
\text{inv-rule-2} & \iff \text{hold}(X), \text{clear}(Y) \Rightarrow \text{get}(X), \text{hold}(Y).
\end{align*}
\]

A run of the inverse program for the previously generated result of ‘hold(cup),
clear(box)’, would produce a single result of ‘get(cup), empty, clear(box)’,
which is not the actual input that was used during the forward run. The inverse
program is now transformed using the proposed exhaustive execution transfor-
mation (but without the pruning rules) to become as follows:

\[
\begin{align*}
\text{start} & \iff \text{depth}(0). \\
\text{reset_rule} & \iff \text{hold}(X) \Rightarrow \text{hold}(X,0,Z1) \iff Z1 < Z \mid \text{hold}(X).
\end{align*}
\]

Running the inverse and transformed program yields the search tree depicted
in Figure 4 (note that the \text{start} constraint is added to trigger the exhaustive
execution). The resultant states in the figure represent all possible inputs to the
forward program and amongst them is the particular input used in Example 1.
Any final state reached by the transformed program is a valid input to a forward
run that generates the goal ‘hold(cup), clear(box)’.
Fig. 4. Derivation tree for Example 2 with query: ‘hold(cup), clear(box), start’

6 Conclusion

A source-to-source transformation was proposed, which expresses any CHR program as one utilizing disjunction, to force an exhaustive explorative execution strategy. It enables a high-level form of execution control that empowers a comprehensive execution while retaining expressive power. It is particularly useful for non-confluent programs with overlapping heads, as it enables finding all possible results to a query. The operational semantics of CHR features a “don’t care non-determinism” where all choices will lead to a successful derivation, so we do not care which one is chosen. The proposed transformation changes it to a “don’t know non-determinism”, where some of the choices will lead to a successful search but we do not know which one beforehand. This change enables exploring the search space generated during any derivation.

The proposed transformation focuses on simplification rules, since propagation and simpagation rules can be represented as simplification rules with a token store history. Therefore, in the future the transformation can be extended to transform all CHR rule types.

The execution platform proposed makes extensive use of the disjunctive operator of CHR $\lor$, and produces a comprehensive search tree for any query. In the future, this transformation can be easily extended with implementations of search strategies, such as the breadth-first transformation [5] due to the presence of disjunction in the rule bodies. The integration and customization of other search strategies can also be incorporated into the transformation.

Despite the pruning rules, the transformation still produces many redundant states (or nodes) which have already been visited in the search tree. Optimizations can be devised to eliminate nodes (or sub-trees) that have already been visited, this could be implemented by the use of a traversal history mechanism.

Furthermore, it would be interesting to add priorities to the branches of the search trees generated and hence enable a priority-based execution. For agent-based programs, this would allow introducing a kind of reasoning whilst performing the various actions.
Bibliography


Extending the 2D DP Framework for CTRSs

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Abstract. Recently, a new dependency pair framework for proving operational termination of Conditional Term Rewriting Systems (CTRSs) has been introduced. We call it 2D DP Framework for CTRSs because it makes explicit and exploits the bidimensional nature of the termination behavior of conditional rewriting: a horizontal component concerning infinite sequences of rewriting steps, and a vertical component that captures infinitely many climbs during the development of a proof tree for a single rewriting step. In this paper we extend the 2D DP Framework for CTRSs with several powerful processors for proving and disproving operational termination that specifically exploit the structure of conditional rules. We provide the first implementation of the 2D DP Framework as part of the termination tool mu-term. Our benchmarks suggest that, with our new processors, the 2D DP Framework is currently the most powerful technique for proving operational termination of CTRSs.

Keywords: Conditional rewriting, dependency pairs, operational termination.

1 Introduction

In [7], the Dependency Pair Framework [3] for Term Rewriting Systems has been extended to prove operational termination [8] of Conditional Term Rewriting Systems (CTRSs). We faithfully capture the bidimensional nature of infinite computations with CTRSs: besides infinite sequences of rewriting steps (a horizontal dimension), there can be infinitely many failing attempts to satisfy the conditions of the rules when a single rewriting step is attempted (a vertical dimension). This twofold origin of infinite computations is captured by two sets of 2D Dependency Pairs (2D DPs). Besides characterizing operational termination of CTRSs in terms of 2D DPs, a 2D DP Framework for mechanizing proofs of operational termination of CTRSs is introduced. A central notion is that of processor, which transforms our termination problems into sets of simpler termination problems which can then be handled independently. This divide and conquer approach is paramount in the (2D) DP Framework (Section 2.1).

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In [7], only four processors were introduced, and three of them were close to well-known processors in the DP Framework: the SCC processor (which permits the use of graph techniques to decompose termination problems), the Reduction Triple Processor (which uses orderings to simplify termination problems) and a shifting processor which just calls the DP Framework for TRSs from the 2D DP Framework when TRSs, rather than CTRSs, are involved. As the benchmarks in this paper show, with those processors we are unable to outperform existing tools, like AProVE [5] and VMTL [11], that prove operational termination of CTRSs by using transformations into TRSs.

Example 1. The operational termination of the CTRS \( \mathcal{R} \) [10, Example 17]:

\[
\begin{align*}
\text{(1)} \quad a \rightarrow h(b) \\
\text{(2)} \quad a \rightarrow h(c) \\
\text{(3)} \quad f(x) \rightarrow y \Leftrightarrow a \rightarrow h(y) \\
\text{(4)} \quad g(x, b) \rightarrow g(f(c), x) \Leftrightarrow f(b) \rightarrow x, x \rightarrow c
\end{align*}
\]

cannot be proved with the processors in [7]. We define new processors to transform and analyze the (satisfaction of the) conditional part of the rules. They detect that \( x \) in (4) is bound only to \( f(b) \) or \( c \) in any rewriting step with (4). This yields a more precise assessment of the role of the rule in the termination behavior and finally leads to a proof of operational termination of \( \mathcal{R} \).

Our contributions in this paper are the following: (1) we refine the calculus of the graph that we use to represent termination problems (Section 3); (2) we adapt Hirokawa and Middeldorp’s subterm criterion for TRSs [6] (Section 4); (3) we define new processors that exploit the simplification of the conditional part of the rules that participate in a given CTRS problem by either removing conditions that unify or by refining other conditions by narrowing (Section 5); (4) we extend and generalize the narrowing processor for TRSs in [3] (Section 6); (5) we introduce a new processor to specifically disprove operational termination of CTRSs (Section 7); then, (6) we describe the implementation of the 2D-DP framework as part of mu-term [2] and (7) provide some benchmarks showing that the 2D DP Framework with the new processors outperforms all existing tools for proving operational termination of CTRSs (Section 8).

2 Preliminaries

We use the standard notations in term rewriting (see, e.g., [9]). In this paper, \( \mathcal{X} \) denotes a countable set of variables and \( \mathcal{F} \) denotes a signature, i.e., a set of function symbols \( \{f, g, \ldots\} \), each with a fixed \textit{arity} given by a mapping \( \text{ar} : \mathcal{F} \rightarrow \mathbb{N} \). The set of terms built from \( \mathcal{F} \) and \( \mathcal{X} \) is \( \mathcal{T}(\mathcal{F}, \mathcal{X}) \). The symbol labeling the root of \( t \) is denoted as root(\( t \)). The set of variables occurring in \( t \) is Var(\( t \)). Terms are viewed as labelled trees in the usual way. Positions \( p, q, \ldots \) are represented by chains of positive natural numbers used to address subterms \( t|_p \) of \( t \). The set of positions of a term \( t \) is Pos(\( t \)). Given \( \Delta \subseteq \mathcal{F}, \text{Pos}_\Delta(\( t \)) \) denotes the set of positions \( p \in \text{Pos}(\( t \)) \) of subterms \( t|_p \) of \( t \) that are rooted by a symbol in \( \Delta \) (i.e., root(\( t|_p \)) \in \Delta \)). A substitution is a mapping from variables into terms which is homomorphically extended to a mapping from terms to terms. A conditional rule
is written $\ell \to r \Leftarrow s_1 \to t_1, \cdots, s_n \to t_n$, where $\ell, r, s_1, t_1, \ldots, s_n, t_n \in \mathcal{T}(\mathcal{F}, \mathcal{X})$ and $\ell \notin \mathcal{X}$. As usual, $\ell$ and $r$ are called the left- and right-hand sides of the rule, and the sequence $s_1 \to t_1, \cdots, s_n \to t_n$ (often abbreviated to $c$) is the conditional part of the rule. We often write $s_i \to t_i \in c$ to refer to the $i$-th atomic condition in $c$ or $s \to t \in c$ if the position of the atomic condition in $c$ does not matter. Rules $\ell \to r \Leftarrow c$ are classified according to the distribution of variables as follows: type 1 (or 1-rules), if $\text{Var}(r) \subseteq \text{Var}(\ell)$; type 2, if $\text{Var}(r) \subseteq \text{Var}(\ell)$; type 3, if $\text{Var}(r) \subseteq \text{Var}(\ell) \cup \text{Var}(c)$; and type 4, if no restriction is given. A CTRS $\mathcal{R}$ is a set of conditional rules; $\mathcal{R}$ is called an $n$-CTRS if it contains only $n$-rules; A 3-CTRS $\mathcal{R}$ is called deterministic if for each rule $\ell \to r \Leftarrow s_1 \to t_1, \ldots, s_n \to t_n$ in $\mathcal{R}$ and each $1 \leq i \leq n$, we have $\text{Var}(s_i) \subseteq \text{Var}(\ell) \cup \bigcup_{j=1}^{i-1} \text{Var}(t_j)$. We write $s \to_R t$ (resp. $s \to_R^* t$) if there is a closed proof tree for $s \to t$ (resp. $s \to^* t$) using the inference system in Figure 1. A proof tree $T$ (for a CTRS $\mathcal{R}$, using the inference system in Figure 1) is closed whenever it is finite and contains no open goals; it is well-formed if it is either an open goal, or a closed proof tree, or a derivation tree of the form $\xrightarrow{T_{i-1}} G \xrightarrow{T_i} \xrightarrow{T_{i+1}} \cdots$ where there is $i, 1 \leq i \leq n$ such that $T_1, \ldots, T_{i-1}$ are closed, $T_i$ is well-formed but not closed, and $T_{i+1}, \ldots, T_n$ are open goals. A proof tree $T$ for $\mathcal{R}$ is a proper prefix of a proof tree $T'$ if there are one or more open goals $G_1, \ldots, G_n$ in $T$ such that $T'$ is obtained from $T$ by replacing each $G_i$ by a derivation tree $T_i$ with root $G_i$. We denote this as $T \subset T'$. An infinite proof tree for $\mathcal{R}$ is an infinite increasing chain of finite trees, that is, a sequence $\{T_i\}_{i \in \mathbb{N}}$ such that for all $i$, $T_i \subset T_{i+1}$. An infinite proof tree is well-formed if it is an ascending chain of well-formed finite proof trees. Intuitively, well-formed trees are the trees that an interpreter would incrementally build when trying to solve one condition at a time from left to right. A CTRS $\mathcal{R}$ is operationally terminating if no infinite well-formed tree for a goal $s \to^* t$ exists.

### 2.1 2D Dependency Pair Framework for CTRSs

In the 2D DP Framework, operational termination of CTRSs is captured using sets $\text{DP}_H(\mathcal{R})$, $\text{DP}_V(\mathcal{R})$, and $\text{DP}_{VH}(\mathcal{R})$ of 2D DPs [7, Section 4]. $\text{DP}_H(\mathcal{R})$ contains new rules capturing infinite computations due to the application of rewriting steps (infinite rewrite sequences); in contrast $\text{DP}_V(\mathcal{R})$ contains rules representing infinite computations due to failing attempts to prove the satisfaction of a condition in a rule; in this respect, $\text{DP}_{VH}(\mathcal{R})$ represents intermediate (but finite) computations between two consecutive uses of a rule in $\text{DP}_V(\mathcal{R})$. 

\[
\begin{align*}
\text{(Ref)} & \quad \frac{t \to^* t}{s \to^* t} \\
\text{(Cong)} & \quad \frac{f(s_1, \ldots, s_i, \ldots, s_k) \to f(s_1, \ldots, s_i, \ldots, s_k)}{s \to^* t} \\
\text{(Tran)} & \quad \frac{s \to u}{s \to^* t} \\
\text{(Repl)} & \quad \frac{\sigma(s_1) \to^* \sigma(t_1) \ldots \sigma(s_n) \to^* \sigma(t_n)}{\sigma(t) \to \sigma(r)}
\end{align*}
\]
Example 2. For $\mathcal{R}$ in Example 1,
\[
\begin{align*}
\text{DP}_H(\mathcal{R}) : & \quad G(x, b) \rightarrow G(f(c), x) \Leftarrow f(b) \rightarrow x, x \rightarrow c \\
\text{DP}_V(\mathcal{R}) : & \quad G(x, b) \rightarrow F(b)
\end{align*}
\]
and $\text{DP}_{V\mathcal{H}}(\mathcal{R}) = \emptyset$. $\text{DP}_H(\mathcal{R})$ is obtained by marking (here just capitalizing the root symbols of) the left- and right-hand sides $g(x, b)$ and $g(f(c), x)$ of (4) and keeping the conditional part. $\text{DP}_V(\mathcal{R})$ is obtained from the left-hand side $g(x, b)$ of (4) and the first condition $f(b)$ of (4) which are marked and combined as a new rule whose conditional part is empty in this case.

Definition 1. Let $\mathcal{P}, \mathcal{Q}, \mathcal{R}$ be CTRSs. A $(\mathcal{P}, \mathcal{Q}, \mathcal{R})$-chain is a sequence $(\alpha_i)_{i \geq 1}$ where $\alpha_i : u_i \rightarrow v_i \Leftarrow c_i \in \mathcal{P}$, together with a substitution $\sigma$ satisfying that, for all $i \geq 1$, (1) $\sigma(s) \rightarrow^*_{\mathcal{R}} \sigma(t)$ for all $s \rightarrow t \in c_i$ and (2) $\sigma(v_i) \rightarrow^*_{\mathcal{R}} \sigma(u_{i+1})$, where $s \rightarrow^*_{\mathcal{R}} t$ if $s = t$ or there is $l \rightarrow r \Leftarrow c \in \mathcal{Q}$ and a substitution $\theta$ such that $s = \theta(l)$, $t = \theta(r)$ and $\theta(u) \rightarrow^*_{\mathcal{R}} \theta(v)$ for all $u \rightarrow v \in c$. We assume $\text{Var}(\alpha_i) \cap \text{Var}(\alpha_j) = \emptyset$, if $i \neq j$ (rename the variables if necessary). A $(\mathcal{P}, \mathcal{Q}, \mathcal{R})$-chain is called minimal if for all $i \geq 1$, $\sigma(v_i)$ is $\mathcal{R}$-operationally terminating.

Rules $\mathcal{Q}$ in $(\mathcal{P}, \mathcal{Q}, \mathcal{R})$ may contribute connections between $\sigma(v_i)$ and $\sigma(v_{i+1})$ using root steps only (compare with the use of $\mathcal{R}$). We often speak of $(\mathcal{P}, \mathcal{Q}, \mathcal{R}, a)$-chains (or $(\mathcal{P}, \mathcal{Q}, \mathcal{R}, m)$-chains), if arbitrary (minimal) chains are considered. A CTRS problem is a tuple $\tau = (\mathcal{P}, \mathcal{Q}, \mathcal{R}, e)$, where $\mathcal{P}$, $\mathcal{Q}$ and $\mathcal{R}$ are CTRSs, and $e \in \{a, m\}$ is a flag. We call $\tau$ finite if there is no infinite $(\mathcal{P}, \mathcal{Q}, \mathcal{R}, e)$-chain; $\tau$ is infinite if $\mathcal{R}$ is operationally nonterminating or there is an infinite $(\mathcal{P}, \mathcal{Q}, \mathcal{R}, e)$-chain. A deterministic 3-CTRS $\mathcal{R}$ is operationally terminating if the two (initial) CTRS problems $\tau_H = (\text{DP}_H(\mathcal{R}), \emptyset, \mathcal{R}, e)$ and $\tau_V = (\text{DP}_V(\mathcal{R}), \text{DP}_{V\mathcal{H}}(\mathcal{R}), \mathcal{R}, e)$ are finite [7, Theorem 2]. A CTRS processor $\mathcal{P}$ maps CTRS problems into sets of CTRS problems. Alternatively, it can also return “no”. $\mathcal{P}$ is sound if for all CTRS problems $\tau$, $\tau$ is finite whenever $\mathcal{P}(\tau) \neq \text{no}$ and all CTRS problems in $\mathcal{P}(\tau)$ are finite. It is complete if for all CTRS problems $\tau$, we have that $\tau$ is infinite whenever $\mathcal{P}(\tau) = \text{no}$ or when $\mathcal{P}(\tau)$ contains an infinite CTRS problem.

Example 3. According to Example 2, for $\mathcal{R}$ in Example 1, $\tau_H = \{(5)\}, \emptyset, \mathcal{R}, m\)$ and $\tau_V = \{(6)\}, \emptyset, \mathcal{R}, m\)$. Some processors apply if $e = m$ only (e.g., Theorems 2 and 7). With $m$ in $\tau_H$ and $\tau_V$ more processors can be potentially used.

In order to prove a deterministic 3-CTRS operationally terminating, we construct two trees whose nodes are labeled with CTRS problems or “yes” or “no”, and whose roots are labeled with $\tau_H$ and $\tau_V$, respectively. For every inner node labeled with $\tau$, there is a processor $\mathcal{P}$ satisfying one of the following conditions: (1) $\mathcal{P}(\tau) = \text{no}$ and the node has just one child that is labeled with “no”; (2) $\mathcal{P}(\tau) = \emptyset$ and the node has just one child that is labeled with “yes”; (3) $\mathcal{P}(\tau) \neq \text{no}$, $\mathcal{P}(\tau) \neq \emptyset$, and the children of the node are labeled with the CTRS problems in $\mathcal{P}(\tau)$. If all leaves of both trees are labeled with “yes” and all used processors are sound, then $\mathcal{R}$ is operationally terminating. If there is a leaf labeled with “no” in one of the trees, and all processors used on the path from the root to this leaf are complete, then $\mathcal{R}$ is operationally nonterminating.
Remark 1. In the following, when defining a processor \( P \) on a CTRS problem \( \tau \), we describe its specific action under the specified conditions. If such conditions do not hold, then we assume \( P(\tau) = \{\tau\} \) (it does nothing). Furthermore, we tacitly assume \( P, Q \) and \( R \) to be deterministic 3-CTRSs.

3 Removing Useless Connection Pairs

Connections between rules in \( P \) become arcs in the graph \( G(\tau) \) associated to \( \tau \), whose set of nodes is \( P \); there is an arc from \( \alpha \) to \( \alpha' \) iff \( (\alpha, \alpha') \) is a \((P, Q, R, e)\)-chain [7, Definition 9]. We estimate them using abstraction and unification [7, Section 6.1], thus obtaining the estimated graph \( EG(\tau) \) [7, Definition 11]. Terms \( s \) and \( t \) unify (written \( s \sim t \)) if there is a substitution \( \sigma \) (called a unifier) such that \( \sigma(s) = \sigma(t) \); when this happens, there is a most general unifier (mgu) \( \theta \) of \( s \) and \( t \) which is unique for \( s \) and \( t \) up to renaming of variables. In the following, if terms \( s \) and \( t \) unify with mgu \( \theta \), we often write \( s \sim^\theta t \) to make it explicit.

Example 4. Consider the following deterministic 3-CTRS \( R \) [7]:

\[
\begin{align*}
g(a) &\rightarrow c(b) & b &\rightarrow f(a) & f(x) &\rightarrow y \Leftarrow g(x) \rightarrow c(y) \\
DP_H(R) : & G(a) \rightarrow B \quad (7) & DP_{VH}(R) : & G(a) \rightarrow B \quad (9) \\
DP_V(R) : & F(x) \rightarrow G(x) \quad (8) & B &\rightarrow F(a) \quad (10)
\end{align*}
\]

For \( \tau_V = \{(8)\}, \{(9), (10)\}, R, m \), the estimated graph \( EG(\tau_V) \) is

\[
\begin{array}{c}
8 \\
\end{array}
\]

The ‘component \( Q \)’ of \( \tau_V \) is essential: the arc is settled because the right-hand side \( G(a) \) of (8) ‘root-rewrites’ into \( F(a) \) using \( DP_{VH}(R) : G(a) \xrightarrow{A(9)} B \xrightarrow{A(10)} F(a) \). No rewriting with \( R \) is possible. Now, \( F(a) \) unifies with the left-hand side \( F(x) \).

In some cases, no pair in \( Q \) can be used for this purpose and we can safely remove \( Q \) from the CTRS problem. This is the focus of our first processor (Theorem 1). The following definitions are necessary. Following [4] (also [7], where we used it to estimate the graph), we let \( TCAP_R \) be \( TCAP_R(x) = y \) if \( x \) is a variable, and

\[
TCAP_R(f(t_1, \ldots, t_k)) = \begin{cases} f(TCAP_R(t_1), \ldots, TCAP_R(t_k)) & \text{if} \ \forall \ell \rightarrow r \Leftarrow c \in R, \\
\ell \text{ and } f(TCAP_R(t_1), \ldots, TCAP_R(t_k)) \text{ does not unify} & y \text{ otherwise} \\
\end{cases}
\]

where \( y \) is a fresh variable. We assume that \( \ell \) shares no variable with \( f(TCAP_R(t_1), \ldots, TCAP_R(t_k)) \). As discussed in [4, 7], with \( TCAP_R \) we approximate reachability problems by means of unification.

Theorem 1 (Removing connection pairs). Let \( P, Q \) and \( R \) be CTRSs. Let \( Q_c = \{ u' \rightarrow v' \Leftarrow c' \in Q \mid \exists u \rightarrow v \Leftarrow c \in P, TCAP_R(v') =^? u \} \). If \( Q_c = \emptyset \), then

\[
P_{RQ}(P, Q, R, c) = \{(P, \emptyset, R, c)\}
\]

is a sound and complete processor.
Example 5. For the CTRS \( \mathcal{R} \) (from the Termination Problem Data Base, TPDB, file TRS\_Conditional/Mixed\_CTRS/fib.xml):

\[
\begin{align*}
\text{plus}(x, y) & \rightarrow y' \iff x \rightarrow 0, y \rightarrow y' \\
\text{fib}(0) & \rightarrow \text{pair}(0, s(0)) \\
\text{fib}(s(x)) & \rightarrow \text{pair}(z, \text{plus}(y, z)) \iff \text{fib}(x) \rightarrow \text{pair}(y, z)
\end{align*}
\]

we have:

\[
\begin{align*}
\text{DP}_H(\mathcal{R}) : & \quad \text{PLUS}(x, y) \rightarrow \text{PLUS}(x', y') \iff x \rightarrow s(x'), y \rightarrow y' \quad (11) \\
\text{FIB}(s(x)) & \rightarrow \text{PLUS}(y, z) \iff \text{fib}(x) \rightarrow \text{pair}(y, z) \quad (12) \\
\text{DP}_V(\mathcal{R}) : & \quad \text{FIB}(s(x)) \rightarrow \text{FIB}(x) \quad (13)
\end{align*}
\]

\( \text{DP}_{VH}(\mathcal{R}) = \text{DP}_H(\mathcal{R}) \), and \( \tau_H = (\{ (11), (12) \}, \emptyset, \mathcal{R}, m) \). For the rhs \( \text{PLUS}(x', y') \) of (11), we have that \( \text{TCAP}_{R, c} \)(PLUS \((x', y')\)) is not unifiable. A similar observation holds for the rhs \( \text{PLUS}(y, z) \) of (12). Thus, for \( \tau_V = (\{ (13) \}, \{ (11), (12) \}, \mathcal{R}, m) \), \( \mathcal{P}_{RQ}(\tau_V) = \{ \tau_{V1} \} \), where \( \tau_{V1} = (\{ (13) \}, \emptyset, \mathcal{R}, m) \).

\( \mathcal{P}_{RQ} \) can be seen as a refinement of \( \mathcal{P}_{\text{SCC}} \) [7], which uses a similar approach.

Remark 2 (Notation). Given (possibly empty) sets of rules \( \mathcal{R}, \mathcal{S} \) and a rule \( \ell \rightarrow r \leftarrow c \), we denote the (possible) replacement of \( \alpha \) in \( \mathcal{R} \) by the rules \( \mathcal{S} \) as:

\[
\mathcal{R}[\mathcal{S}]_\alpha = \begin{cases} (\mathcal{R} - \{ \alpha \}) \cup \mathcal{S} & \text{if } \alpha \in \mathcal{R} \\ \mathcal{R} & \text{otherwise} \end{cases}
\]

We let \( c[c_i'] \) be the condition obtained by replacing in \( c \) the \( i \)-th atomic condition \( s_i \rightarrow t_i \in c \) by the sequence of conditions \( c' \). If \( c' \) is empty, we write \( c[\emptyset]_i \).

4 Subterm processor

In this section, we generalize the subterm processor for TRSs [6]. In the following, we write \( s \succeq t \) iff \( t \) is a subterm of \( s \), and \( s \succeq t \) iff \( s \succeq t \) and \( s \neq t \).

Definition 2 (Root symbols of a CTRS). Let \( \mathcal{R} \) be a CTRS. The set of root symbols associated to \( \mathcal{R} \) is:

\[
\text{Root}(\mathcal{R}) = \{ \text{root}(\ell) \mid \ell \rightarrow r \leftarrow c \in \mathcal{R} \} \cup \{ \text{root}(r) \mid \ell \rightarrow r \leftarrow c \in \mathcal{R}, r \notin X \}
\]

Definition 3 (Simple projection). Let \( \mathcal{R} \) be a CTRS. A simple projection for \( \mathcal{R} \) is a mapping \( \pi : \text{Root}(\mathcal{R}) \rightarrow \mathbb{N} \) such that \( \pi(f) \in \{ 1, \ldots, \text{ar}(f) \} \). The mapping that assigns a subterm \( \pi(t) = t|_{\pi(f)} \) to each term \( t \) with root \( \ell \) in \( \text{Root}(\mathcal{R}) \) is also denoted by \( \pi \); we also let \( \pi(x) = x \) if \( x \in X \).

Given a simple projection \( \pi \) for a CTRS \( \mathcal{R} \), we let \( \pi(\mathcal{R}) = \{ \pi(\ell) \rightarrow \pi(r) \mid \ell \rightarrow r \leftarrow c \in \mathcal{R} \} \). Note that the conditions are just dismissed. Given a CTRS problem \( (\mathcal{P}, \mathcal{Q}, \mathcal{R}, c) \), the subterm processor removes from \( \mathcal{P} \) and \( \mathcal{Q} \) those rules \( u \rightarrow v \leftarrow c \) whose left-hand side \( u \) contains an immediate subterm \( \pi(u) \) which is a strict superterm of an immediate subterm \( \pi(v) \) of \( v \) (i.e., \( \pi(u) \triangleright \pi(v) \)). In the following result we say that a CTRS \( \mathcal{R} \) is collapsing if there is a rule \( \ell \rightarrow r \leftarrow c \in \mathcal{R} \) such that \( r \) is a variable. We also write \( \mathcal{D}_\mathcal{R} \) to denote the set of symbols which are defined by rules in \( \mathcal{R} \), i.e., \( \mathcal{D}_\mathcal{R} = \{ \text{root}(\ell) \mid \ell \rightarrow r \leftarrow c \in \mathcal{R} \} \).
Theorem 2 (Subterm processor). Let $P$, $Q$, and $R$ be CTRSs such that $P$ and $Q$ are not collapsing and $(\text{Root}(P) \cup \text{Root}(Q)) \cap D_R = \emptyset$. Let $\pi$ be a simple projection for $P \cup Q$ and $\alpha : u \rightarrow v \in P \cup Q$ be such that $\pi(P) \cup \pi(Q) \subseteq \triangleright$ and $\pi(u) \triangleright \pi(v)$. Then,

$$P_{\triangleright}(P, Q, R, m) = \{(P[\emptyset]_\alpha, Q[\emptyset]_\alpha, R, m)\}$$

is a sound and complete processor.

Example 6. For $\tau_{V1} = (\{(13)\}, \emptyset, R, m)$ in Example 5, with $\pi(\text{FIB}) = 1$, $P_{\triangleright}(\tau_{V1}) = \{(\emptyset, \emptyset, R, m)\}$ because $\pi(\text{FIB}(s(x))) = s(x) \triangleright x = \pi(\text{FIB}(x))$. Since $(\emptyset, \emptyset, R, m)$ is trivially finite, $\tau_{V1}$ is also finite.

$P_{\triangleright}$ does not apply to $\tau_H$ in Example 5 to remove (12). This is because no variable in the right-hand side $\text{PLUS}(y, z)$ is in the corresponding left-hand side $\text{FIB}(s(x))$. Thus $s(x)$ cannot be seen as a superterm of any subterm of $\text{PLUS}(y, z)$.

5 Simplifying the Conditions of the Rules

The condition $c$ of a rule $\ell \rightarrow r \leftarrow c$ controls its applicability to specific redexes $\sigma(\ell)$ depending on the satisfaction of the instantiated condition $\sigma(c)$. However, dealing with the conditional part of the rules often requires specialized techniques (for instance, for solving conditional constraints as in the removal triple processor of [7, Theorem 11]) which can make proofs difficult. In other cases, the rules are just useless, because they cannot be applied at all, but establishing uselessness to take advantage of it in the analysis of termination is often difficult.

Example 7. Consider the CTRS $R$ [9, Example 7.2.45]:

\begin{align*}
a &\rightarrow a \leftarrow b \rightarrow x, \ c \rightarrow x & (14) \\
b &\rightarrow d \leftarrow d \rightarrow x, \ e \rightarrow x & (15)
\end{align*}

where $a, \ldots, e$ are constants and $x$ is a variable. Powerful tools like AProVE do not find a proof of operational termination of $R$ by using transformations. Our implementation of the processors in [7] cannot prove it either. Clearly, (15) and (16) cannot be used in any rewriting step (in any $(P, Q, R)$-chain) because $d$ and $e$ are irreducible and the only way to get the condition $d \rightarrow x, e \rightarrow x$ satisfied is the instantiation of $x$ to both $d$ and $e$, which is not possible. The only processor in [7] which specifically addresses the use of the information in the conditions of the rules is $P_{UR}$ which removes unsatisfiable rules [7, Theorem 11], i.e., those rules which cannot be used in any chain because the conditional part is unsatisfiable. $P_{UR}$ uses a well-founded ordering $\sqsubseteq$ which, under appropriate conditions, removes $\ell \rightarrow r \leftarrow c$, if $\sigma(t) \sqsubseteq \sigma(s)$ holds for some $s \rightarrow t \in c$ and all substitutions $\sigma$. In our example, we should have $x \sqsubseteq d$ or $x \sqsubseteq e$ (for all instances of $x$). This is clearly impossible because, in particular, it would require $d \sqsubseteq d$ or $e \sqsubseteq e$, which contradicts well-foundedness of $\sqsubseteq$. Some preprocessing on the conditions is required before achieving a proof (see Examples 9 and 10 below).

The common feature of the new processors in this section is their simplification of the conditional part $c$ of the rules in the CTRSs of a given CTRS problem.
5.1 Simplification by unification

Some conditions $s \rightarrow t$ in the conditional part $c$ of a rule $\ell \rightarrow r \Leftarrow c$ cannot start any (specific) rewriting computation before reaching the corresponding instance of $t$. They can then be viewed as unification problems $s \equiv t$. Therefore, $s \rightarrow t$ can be removed from $c$ if we instantiate the rule with the most general unifier $\theta$ of $s$ and $t$. In the following, given a CTRS $R$, we say that a non-variable term $t$ is a narrowing redex (or a narrex, for short) if there is a rule $\ell \rightarrow r \Leftarrow c \in R$ such that $\text{Var}(t) \cap \text{Var}(\ell) = \emptyset$ and $t$ and $\ell$ unify. We say that $\text{NARR}_R(t)$ holds if $t$ contains a narrex, i.e., there is a non-variable position $p \in \text{Pos}_R(t)$ such that $t|_p$ is an $R$-narrex. In the following results, given a rule $\alpha : \ell \rightarrow r \Leftarrow c$ with $n$ conditions, some $i$, $1 \leq i \leq n$, and a substitution $\theta$, we let $\alpha_{\theta,i}$ be the rule $\alpha_{\theta,i} : \theta(\ell) \rightarrow \theta(r) \Leftarrow \theta(c|_i)$.

**Theorem 3 (Simplifying unifiable conditions).** Let $P$, $Q$, and $R$ be CTRSs. Let $\alpha : \ell \rightarrow r \Leftarrow c \in P \cup Q \cup R$ and $s_i \rightarrow t_i \in c$ such that: (1) $s_i$ is linear, (2) $\text{NARR}_R(s_i)$ does not hold, (3) $s_i =_\theta t_i$, (4) for all $s \rightarrow t \in c|_i$,

\[
\text{Var}(s_i) \cap \text{Var}(s) = \emptyset,
\]

(5) for all $s \rightarrow t \in c$, $\text{Var}(s_i) \cap \text{Var}(t) = \emptyset$, and (6) $\text{Var}(s_i) \cap \text{Var}(r) = \emptyset$. Then,

\[
P_{\text{SUC}}(P, Q, R, e) = (\{P[\{\alpha_{\theta,i}\}]_\alpha, Q[\{\alpha_{\theta,i}\}]_\alpha, R[\{\alpha_{\theta,i}\}]_\alpha, e\})
\]

is a sound and complete processor.

**Remark 3.** Requirement (5) in Theorem 3 (plus determinism of $P$, $Q$, and $R$) implies that $\text{Var}(s_i) \subseteq \text{Var}(\ell)$. Indeed, every variable in $s$ which is not in $\ell$ must occur in some of the (previous) $t$, which is forbidden by (5).

**Example 8.** For $\tau_H$ in Example 5, the estimated graph $\text{EG}(\tau_H)$ is

\[
\begin{array}{c}
11 \\
\hline
12
\end{array}
\]

With $P_{\text{SCC}}$, which decomposes a CTRS problem $\tau$ into as many CTRS problems as Strongly Connected Components (i.e., maximal cycles) are in $\text{EG}(\tau)$ [7, Theorem 9], $P_{\text{SCC}}(\tau_H) = \{\tau_{H1}\}$, where $\tau_{H1} = (\{(11)\}, \emptyset, R, m)$. We use $P_{\text{SUC}}$ to transform $\tau_{H1}$ into $\{\tau_{H2}\}$, where $\tau_{H2} = (\{(17)\}, \emptyset, R, m)$ for the rule

\[
\text{PLUS}(s(x'), y) \rightarrow \text{PLUS}(x', y') \Leftarrow y \rightarrow y'
\]

(17)

Now, $P_{\text{SUC}}(\tau_{H2}) = \{\tau_{H3}\}$, for $\tau_{H3} = (\{(18)\}, \emptyset, R, m)$ with

\[
\text{PLUS}(s(x'), y') \rightarrow \text{PLUS}(x', y')
\]

(18)

$P_{\sigma}$ removes (18) from $\tau_{H3}$ (with $\pi(\text{PLUS}) = 1$), to yield a trivially finite CTRS problem which proves $\tau_H$ in Example 5 finite. With the proof of finiteness of $\tau_V$ in Example 6, we conclude that $R$ in Example 5 is operationally terminating.
Example 9. For $\mathcal{R}$ in Example 7, we have:

$$\begin{align*}
DP_H(\mathcal{R}) : A &\rightarrow A \leftarrow b \rightarrow x, c \rightarrow x \quad (19) \\
DP_V(\mathcal{R}) : A &\rightarrow B \quad (20) \\
DP_{VH}(\mathcal{R}) = \emptyset, \quad \tau_H = (DP_H(\mathcal{R}), \emptyset, \mathcal{R}, m), \text{ and } \tau_V = (DP_V(\mathcal{R}), DP_{VH}(\mathcal{R}), \mathcal{R}, m).
\end{align*}$$

We use $P_{SUC}$ (twice) to simplify (15) and (16) to transform $\tau_H$ into $\{\tau_{H1}\}$, where $\tau_{H1} = (\{(19)\}, \emptyset, \{(14),(22),(23)\}, m)$ with

$$b \rightarrow d \leftarrow e \rightarrow d \quad (22) \quad c \rightarrow d \leftarrow e \rightarrow d \quad (23)$$

Now, in contrast to $P_{UR}$ in [7], our following processor uses a syntactic criterion to remove those rules that cannot be used due to the unsatisfiability of the conditional part of the rule.

**Theorem 4 (Irreducible conditions).** Let $\mathcal{P}$, $\mathcal{Q}$, and $\mathcal{R}$ be CTRSs. Let $\alpha : \ell \rightarrow r \leftarrow c \in \mathcal{P} \cup \mathcal{Q} \cup \mathcal{R}$ and $s \rightarrow t \in c$ be such that: (1) $s$ is linear, (2) $\text{Narr}_\mathcal{R}(s)$ does not hold, and (3) $s$ and $t$ do not unify. Then,

$$P_{IC}(\mathcal{P}, \mathcal{Q}, \mathcal{R}, e) = \{(P[\sigma]_{\alpha}, Q[\sigma]_{\alpha}, R[\sigma]_{\alpha}, e)\}$$

is a sound and (if $\alpha \notin \mathcal{R}$ or $e = a$) complete processor.

Example 10. For $\tau_{H1}$ in Ex. 9, $P_{IC}(\tau_{H1}) = \{\tau_{H2}\}$, with $\tau_{H2} = (\{(19)\}, \emptyset, \{(14)\}, m)$. With $P_{SUC}$, we obtain $\tau_{H3} = (\{(24)\}, \emptyset, \{(14)\}, m)$, where

$$A \rightarrow A \leftarrow c \rightarrow b \quad (24)$$

Note that this could not be done directly on $\tau_H$ because $\text{Narr}_\mathcal{R}(b)$ is true. In contrast, $\text{Narr}_\{(14)\}(b)$ is not. Now, $P_{IC}(\tau_{H3}) = \{(\emptyset, \emptyset, \{(14)\}, m)\}$, which is finite, thus proving $\tau_H$ in Example 9 finite as well. We prove $\tau_V$ finite by using $P_{SCC}$ in [7] and then conclude operational termination of $\mathcal{R}$ in Example 7.

### 5.2 Narrowing the conditions of the rules

Reachability problems $\sigma(s) \rightarrow^* \sigma(t)$ are often investigated using narrowing and unification conditions directly over terms $s$ and $t$, thus avoiding the ‘generation’ of the required substitution $\sigma$. In the following we define the notion of narrowing that we use for CTRSs as a suitable extension of the usual definition for TRSs.

**Definition 4 (Narrowing with CTRSs).** Let $\mathcal{R}$ be a CTRS. A term $s$ narrows to a term $t$ (written $s \overset{\ell}{\rightarrow}_{\mathcal{R}, \theta} t$ or just $s \overset{\ell}{\rightarrow} t$), iff there is a nonvariable position $p \in \text{Pos}_\mathcal{F}(s)$ and a rule $\ell \rightarrow r \leftarrow \bigwedge_{i=1}^{n} s_i \rightarrow t_i$ in $\mathcal{R}$ (sharing no variable with $s$) such that:

1. $s|_p \overset{\theta}{=} \ell$,
2. for all $i$, $1 \leq i \leq n$, $\vartheta_{i-1}(s_i) \overset{\ell_i}{\rightarrow}_{\mathcal{R}, \theta, \vartheta_i} t'_i$ and $t'_i \overset{\vartheta_i}{=} \vartheta_i(\vartheta_{i-1}(t_i))$, where $\vartheta_0 = \theta_0$ and for all $i > 0$, $\vartheta_i = \vartheta_i \circ \vartheta_{i-1}$, and
3. $t = \theta(s[p], \vartheta_t)$, where $\theta = \vartheta_n$. 


The reflexive and transitive closure $\sim^*$ of $\sim$ is $\sim^* = \bigcup_{i \geq 0} \sim^i$, where $s \sim_0^* s$, and $s \sim_n^* t$ if $s \sim_{\mathcal{R},\theta_n} u$, $u \sim_{\mathcal{R},\theta_{n-1}}^* t$, and $\theta_n = \theta_{n-1} \circ \theta$. In all narrowing steps we assume that a renamed rewrite rule $\ell \rightarrow r \leftarrow c$ has been used in such a way that no variable in the rule occurs in any previous term in the sequence.

We say that a CTRS $\mathcal{R}$ has no strict overlaps if for all $\alpha : \ell \rightarrow r \leftarrow c \in \mathcal{R}$ and $p \in \text{Pos}_\mathcal{R}(\ell) - \{\ell\}$, there is no rule $\alpha' : \ell' \rightarrow r' \leftarrow c'$ such that $\text{Var}(\ell) \cap \text{Var}(\ell') \neq \emptyset$ (rename the variables if necessary; $\alpha$ and $\alpha'$ can be the same rule) $\ell|_p$ and $\ell'$ unify. In the following results, given a CTRS $\mathcal{R}$, we use the following notation:

- $\text{Rules}(\mathcal{R}, f) = \{l \rightarrow r \leftarrow c \in \mathcal{R} \mid \text{root}(l) = f\}$ is the set of rules in $\mathcal{R}$ defining a symbol $f$ and $\text{Rules}(\mathcal{R}, t) = \cup_{f \in \mathcal{F}(t)} \text{Rules}(\mathcal{R}, f)$ is the set of rules in $\mathcal{R}$ defining the symbols in term $t$.
- Given a term $s$, $N_1(\mathcal{R}, s)$ represents the set of one-step $\mathcal{R}$-narrowings issued from $s$: $N_1(\mathcal{R}, s) = \{(t,c,\theta) \mid s \sim_{\ell \rightarrow r \leftarrow c,\theta} t, \ell \rightarrow r \leftarrow c \in \mathcal{R}\}$.
- Given a rule $\alpha : \ell \rightarrow r \leftarrow c$ with $n$ conditions and $i, 1 \leq i \leq n$, we let $N(\mathcal{R}, \alpha, i) = \{\theta(\ell) \rightarrow \theta(r) \leftarrow \theta(c)[\theta(c'), w \rightarrow \theta(t_i)], s_i 
\rightarrow t_i \in c, (w, c', \theta) \in N_1(\mathcal{R}, s_i)\}$

**Theorem 5 (Narrowing the conditions of rules).** Let $\mathcal{P}$, $\mathcal{Q}$, and $\mathcal{R}$ be CTRSs. Let $\alpha : \ell \rightarrow v \leftarrow c \in \mathcal{P} \cup \mathcal{Q}$ and $s_i \rightarrow t_i \in c$ be such that: (1) $s_i$ is linear, (2) $\text{Rules}(\mathcal{R}, s_i)$ have no strict overlap, (3) $\text{Var}(s_i) \cap \text{Var}(t_i) = \emptyset$, and (4) $s_i$ and $t_i$ do not unify. Then,

$$\mathcal{P}_{\text{NC}}(\mathcal{P}, \mathcal{Q}, \mathcal{R}, e) = \{(\mathcal{P}[N(\mathcal{R}, \alpha, i)]|_{\alpha}, \mathcal{Q}[N(\mathcal{R}, \alpha, i)]|_{\alpha}, \mathcal{R}, e)\}$$

is a sound and complete processor.

The unification requirement is essential for the correctness of $\mathcal{P}_{\text{NC}}$.

**Example 11.** Consider the CTRS $\{a \rightarrow b, c \rightarrow d \leftarrow a \rightarrow a\}$. The left-hand side $a$ of the condition in the second rule narrows into $b$. But $c \rightarrow d \leftarrow b \rightarrow a$ (as obtained by $\mathcal{P}_{\text{NC}}$), now forbids the rewriting step $c \rightarrow d$.

In order to avoid the problem illustrated by the previous example, our next processor is able to remove the requirement of unification by combining the transformations in Theorems 3 and 5 into a single processor.

**Theorem 6 (Simplification and narrowing).** Let $\mathcal{P}$, $\mathcal{Q}$, and $\mathcal{R}$ be CTRSs. Let $\alpha : \ell \rightarrow v \leftarrow c \in \mathcal{P} \cup \mathcal{Q}$ and $s_i \rightarrow t_i \in c$ be such that: (1) $s_i$ is linear, (2) $\text{Rules}(\mathcal{R}, s_i)$ have no strict overlap, (3) $\text{Var}(s_i) \cap \text{Var}(t_i) = \emptyset$, (4) $s_i =^{?} t_i$, (5) for all $s_j \rightarrow t_j \in c, i \neq j, \text{Var}(s_i) \cap \text{Var}(s_j) = \emptyset$, (6) for all $s_j \rightarrow t_j \in c, \text{Var}(s_i) \cap \text{Var}(t_j) = \emptyset$, and (7) $\text{Var}(s_i) \cap \text{Var}(w) = \emptyset$. Then,

$$\mathcal{P}_{\text{SUNC}}(\mathcal{P}, \mathcal{Q}, \mathcal{R}, e) = \{(\mathcal{P}[N(\mathcal{R}, \alpha, i) \cup \{\alpha_{\theta,i}\}]|_{\alpha}, \mathcal{Q}[N(\mathcal{R}, \alpha, i) \cup \{\alpha_{\theta,i}\}]|_{\alpha}, \mathcal{R}, e)\}$$

is a sound and complete processor.
Example 12. For \( \tau_H \) in Example 3, rule (5) is transformed by \( P_{\text{SUNC}}(\tau_H) = \{ \tau_{H1} \} \), with \( \tau_{H1} = \{(25),(26),(27)\}, \emptyset, R, m \) where

\[
\begin{align*}
G(x,b) & \rightarrow G(f(c),x) \leftarrow a \rightarrow h(b), b \rightarrow x, x \rightarrow c & (25) \\
G(x,b) & \rightarrow G(f(c),x) \leftarrow a \rightarrow h(c), c \rightarrow x, x \rightarrow c & (26) \\
G(f(b),b) & \rightarrow G(f(c),f(b)) \leftarrow f(b) \rightarrow c & (27)
\end{align*}
\]

Using now \( P_{\text{SUC}} \) twice we obtain \( \tau_{H2} = \{(28),(29),(27)\}, \emptyset, R, m \), where

\[
\begin{align*}
G(b,b) & \rightarrow G(f(c),b) \leftarrow a \rightarrow h(b), b \rightarrow c & (28) \\
G(c,b) & \rightarrow G(f(c),c) \leftarrow a \rightarrow h(c), c \rightarrow c & (29)
\end{align*}
\]

We use \( P_{\text{IC}} \) to remove (28) from \( \tau_{H2} \) due to its condition’s unsatisfiability. The graph for the obtained problem \( \tau_{H3} = \{(27),(29)\}, \emptyset, R, m \) is

\[
\begin{array}{c}
27 \\
\rightarrow
\end{array}
\]

With \( P_{\text{SCC}} \) we obtain \( \tau_{H4} = P_{\text{SCC}}(\tau_{H3}) = \{(27)\}, \emptyset, R, m \). The proof continues by proving \( \tau_{H4} \) finite using the sequence \( P_{\text{NC}}, P_{\text{IC}}, P_{\text{SUC}}, \) and \( P_{\text{NR}} \) (see Theorem 7 below), with intermediate applications of \( P_{\text{SCC}} \). The complete proof can be found in the benchmarks page (see below) under the label \text{jlap09-ex17.trs}.

6 Narrowing the Right-hand Sides of Rules

As mentioned in Section 5.2, reachability problems are often approximated or advanced by narrowing. The connection between two rules \( \alpha : u \rightarrow v \leftarrow c \) and \( \alpha' : u' \rightarrow v' \leftarrow c' \in \mathcal{P} \) within a \((\mathcal{P},\mathcal{Q},\mathcal{R})\)-chain is a kind of reachability problem \( \sigma(v)(\rightarrow_{\mathcal{R}}^n \circ \Delta_{\alpha} \circ \sigma(u')) \), which can also be investigated using narrowing. If there is a (nontrivial, i.e., involving some rewriting with \( \mathcal{R} \) or \( \mathcal{Q} \)) connection between \( \alpha \) and \( \alpha' \) as above, then after narrowing \( v \) into all its possible narrowings \( v_1,\ldots,v_n \), the connection will be exhibited by some of the \( v_i \). The good point is that connections between \( v \) and \( u \) that are obtained by the approximations, but which are unfeasible, may be removed by other processors (typically \( P_{\text{SCC}} \)), thus leading to a more precise analysis. In the following, given a CTRS \( \mathcal{R} \), a rule \( \alpha : u \rightarrow v \leftarrow c \) and a narrowing step \( v \sim_{r \rightarrow r \leftarrow c,\theta} w \) on the right-hand side \( v \) of \( \alpha \), we say that \( \theta(u) \rightarrow w \leftarrow \theta(c), \theta(d) \) is a narrowing of \( \alpha \). Thus, we let

\[
\mathcal{N}(\mathcal{R},\alpha) = \{ \theta(u) \rightarrow w \leftarrow \theta(c), \theta(d) | v \sim_{r \rightarrow r \leftarrow c,\theta} w, \alpha = u \rightarrow v \leftarrow c, \ell \rightarrow r \leftarrow d \in \mathcal{R} \}
\]

Theorem 7 (Narrowing with \( \mathcal{R} \)). Let \( \mathcal{P}, \mathcal{Q}, \) and \( \mathcal{R} \) be CTRSs. Let \( u \rightarrow v \leftarrow c \in \mathcal{P} \) be such that \( v \) is linear, and for all \( u' \rightarrow v' \leftarrow c' \in \mathcal{P} \cup \mathcal{Q} \) (with possibly renamed variables), \( v \) and \( u' \) do not unify. Then,

\[
P_{\text{NR}}(\mathcal{P},\mathcal{Q},\mathcal{R},m) = \{ (\mathcal{P}|\mathcal{N}(\mathcal{R},\alpha))_\alpha, \mathcal{Q}, \mathcal{R}, m \}
\]

is a sound and complete processor.
In the following processor, we use the rules in $Q$ to narrow the right-hand sides of pairs $u \rightarrow v \Leftarrow c \in P$ at the root only. We now let

$$\mathcal{N}_d(R,\alpha) = \{\theta(u) \rightarrow w \Leftarrow \theta(c),\theta(d) \mid v \xi^{\ell \rightarrow r \Leftarrow d,\theta} w,\alpha = u \rightarrow v \Leftarrow c,\ell \rightarrow r \Leftarrow d \in R\}$$

**Theorem 8 (Narrowing with $Q$).** Let $P$, $Q$, and $R$ be CTRSs. Let $u \rightarrow v \Leftarrow c \in P$ be such that $v$ is linear, $\text{NARR}_R(v)$ does not hold, and for all $u' \rightarrow v' \Leftarrow c' \in P$ (with possibly renamed variables), $v$ and $u'$ do not unify. Then,

$$P_{NQ}(P,Q,R,e) = \{(P[\mathcal{N}_d(Q,\alpha)]\alpha,Q,R,a)\}$$

is a sound and complete processor.

**Example 13.** We apply $P_{NQ}$ to $\tau_V$ in Example 4 to obtain $P_{NQ}(\tau_V) = \{\tau_{V1}\}$ where $\tau_{V1} = (\{(30)\},\{(9),(10)\},R,m)$ with

$$F(a) \rightarrow B$$

(30)

And yet $P_{NQ}(\tau_{V1}) = \{\tau_{V2}\}$ where $\tau_{V2} = (\{(31)\},\{(9),(10)\},R,a)$ with

$$F(a) \rightarrow F(a)$$

(31)

This is an infinite CTRS problem, which we will handle with our last processor, introduced in the following section.

## 7 Detection of Infinite CTRS Problems

The following processor detects a simple kind of infinite CTRS problems.

**Theorem 9 (Infinite problem).** Let $P$, $Q$, and $R$ be CTRSs. Let $u \rightarrow v \Leftarrow c \in P$ and $\vartheta,\vartheta'$ be substitutions such that for all $s \rightarrow t \in c$, $\vartheta'(s) \rightarrow^{*}_R \vartheta'(t)$ and $\vartheta(v) = \vartheta(\vartheta'(u))$. Then,

$$P_{Inf}(P,Q,R,e) = \text{no}$$

is a sound and complete processor.

Note that ensuring the existence of a substitution $\vartheta$ that guarantees that the reachability conditions $\vartheta(s) \rightarrow^{*}_R \vartheta(t)$ hold is essential.

**Example 14.** With $\tau_{V2}$ in Example 4 we have $P_{Inf}(\tau_{V2}) = \text{no}$, witnessing that $R$ in Example 4 is not operationally terminating.

In Example 14 we easily proved $\tau_{V2}$ in Example 4 infinite because rule (31) has no conditional part. Thus, we do not need to check the condition $\vartheta(s) \rightarrow^{*}_R \vartheta(t)$ prescribed in Theorem 9, and we actually let $\vartheta$ be the identity substitution. For pairs $u \rightarrow v \Leftarrow c \in P$ where $c$ is not empty, we can use the processor only if we find a substitution $\vartheta$ such that $\vartheta(s) \rightarrow^{*}_R \vartheta(t)$ holds for all $s \rightarrow t \in c$. In general, this is not computable. In our current implementation, we apply $P_{Inf}$ using a rule $u \rightarrow v \Leftarrow c$ only if there is a substitution $\vartheta$ such that $\vartheta(s) = \vartheta(t)$ for all $s \rightarrow t \in c$ and then we check whether $\vartheta(v) = \vartheta(\vartheta(u))$ for some substitution $\theta$. 
8 Experimental Evaluation

This is the first implementation of the 2D-DP framework presented in [7]. It has been developed as part of the tool MU-TERM [2]. In order to assess the practical contributions of the results in this paper, we implemented two versions:

- **MU-TERM 5.11** (WRLA’14), which includes the framework and processors defined in [7].
- **MU-TERM 5.12** (LOPSTR’14), which extends the previous version with the processors defined in this paper.

Note that our processor $P_{inf}$ is the first technique for proving operational non-termination of CTRSs implemented in any termination tool\(^3\). We compared the MU-TERM implementations and the last version of the existing tools that handle conditional rewriting problems: VMTL 1.3 (http://www.logic.at/vmtl/) and AProVE 2014 (http://aprove.informatik.rwth-aachen.de). The experiments have been performed on an Intel Core 2 Duo at 2.4GHz with 8GB of RAM, running OS X 10.9.1 using a 60 seconds timeout. We considered examples from different sources: the CTRSs in the termination problem database\(^4\), TPDB 8.0.7; the CTRSs in the VMTL webpage\(^5\); and the new CTRSs presented in [7] and in this paper. Results are summarized as follows\(^6\):

<table>
<thead>
<tr>
<th>Tool Version</th>
<th>Proved (YES/NO)</th>
<th>Av. YES</th>
<th>Av. NO</th>
</tr>
</thead>
<tbody>
<tr>
<td><strong>MU-TERM 5.11</strong> (WRLA’14)</td>
<td>15/33 (15/0)</td>
<td>0.35s</td>
<td>0s</td>
</tr>
<tr>
<td><strong>MU-TERM 5.12</strong> (LOPSTR’14)</td>
<td>26/33 (21/5)</td>
<td>1.95s</td>
<td>1.02s</td>
</tr>
<tr>
<td>AProVE 2014</td>
<td>18/33 (18/0)</td>
<td>9.20s</td>
<td>0s</td>
</tr>
<tr>
<td>VMTL 1.3</td>
<td>18/33 (14/4)**</td>
<td>6.94s</td>
<td>1.05s*</td>
</tr>
</tbody>
</table>

The practical improvements revealed by the experimental evaluation are promising. We can prove (now) termination of 26 of the 33 examples, 11 more examples than our previous version, including 3 examples that cannot be proved by any other automatic tool. Furthermore, if we consider the 15 problems that can be proved by both MU-TERM versions, the new processors yield a faster tool, witnessed by a speedup of 1.26 with respect to the WRLA version. The new processors are very useful to simplify conditions on rules and are also very helpful to detect non-terminating chains. But there is still room for improvement, as we can see in the results obtained by AProVE and VMTL, where other processors based on transformations (Instantiation, Forward Instantiation, Rewriting, . . . ) and polynomial interpretations using negative numbers can be useful to improve the results obtained in this paper.

\(^3\) With regard to the negative proofs reported here for VMTL, the VMTL web site says that “non-termination means non-termination of the transformed TRS obtained by the conditional TRS through a transformation”. Since the transformation used by VMTL is not complete (see [11]), such negative results do not imply operational nontermination of the original CTRS.

\(^4\) See http://www.termination-portal.org/wiki/TPDB/

\(^5\) See http://www.logic.at/vmtl/benchmarks-cond.html

\(^6\) detailed benchmarks can be found in http://zenon.dsic.upv.es/muterm/benchmarks/lopstr14/benchmarks.html
9 Related work and conclusions

In [7], the generalization of the DP Framework to CTRSs is accomplished, including the definition of sets of dependency pairs which can be used to provide an independent description of infinite computations in the two (horizontal and vertical) dimensions. In particular, the 2D DP Framework was defined, but only a few processors were presented. Furthermore, such processors barely exploit the peculiarities of the 2D DP framework (e.g., the use of an additional set of pairs Q for connecting pairs in P, or considering the conditional part of the rules to remove and transform them in P, Q, and R). Finally, no implementation or experimental analysis of the 2D DP framework was available.

In this paper we have defined 8 new processors: P_{RQ}, which removes pairs from Q which are unable to establish any connection within a (P, Q, R)-chain, P_{P}, which removes pairs from P and Q without paying attention to the structure of rules in R, P_{SUC}, which faithfully removes unifiable conditions, P_{IC}, which removes rules containing conditions that cannot be used in any computation, P_{NC}, which transforms the conditional part of the rules by narrowing, P_{NR}, which transforms the right-hand sides of the rules in P by using narrowing with R, P_{NQ}, which narrows with Q instead, and P_{Inf}, which provides a simple way to detect infinite CTRS problems. We have implemented all our processors (including the ones introduced in [7]). Among these processors, only P_{P}, P_{NR} and P_{Inf} have some analogue processor in the DP Framework for TRSs. All other processors work on the conditional part of the rules or (as P_{NQ}) use active components (e.g., the rules in Q) which are missing in the DP Framework.

Remark 4. Although CTRSs can be transformed into TRSs by means of operational-nontermination-preserving transformations (like U in [9, Definition 7.2.48]), and then existing (narrowing, instantiation) processors of the DP Framework for TRSs (see [3]) could be applied to hopefully obtain similar effects in proofs of operational termination via termination of TRSs, it is unclear (but certainly interesting subject of further research) whether some of our more specific processors P_{SUC}, P_{IC}, P_{SIC}, and P_{NC} could be simulated by some of those processors after the transformation. Our present view is that this is unlikely in most cases. First, it depends on the considered transformation. But even using (a variant of) U as done by AProVE and VMTL, if this connection were easy, our benchmarks would be closer to those of tools like AProVE which implements most processors and techniques of the DP Framework of TRSs. In contrast, the 2014 edition of the International Termination Competition\(^7\) confirms our benchmarks, even with a time-out of five minutes, instead of the one minute time-out used in previous competitions and adopted in our benchmarks, see http://nfa.imn.htwk-leipzig.de/termcomp/show_job_results/5382

for the results of the TRS Conditional subcategory. This suggests that the use of native processors in the 2D DP Framework is better than transforming CTRS problems into TRS problems and then using the DP Framework for TRSs.

This paper gives strong evidence suggesting that the 2D DP Framework is currently the most powerful technique for proving operational termination of CTRSs. For instance, the CTRSs \( R \) in Examples 1 and 7 cannot be proved operationally terminating by AProVE or VMTL. And \( R \) in Example 4 cannot be proved operationally nonterminating because the transformation used by AProVE and VMTL is not complete and does not capture operational nontermination. Furthermore, neither these three examples nor \( R \) in Example 5 can be proved operationally (non)terminating with the processors in [7]. It is true, however, that some of the examples are proved operationally terminating by AProVE or VMTL whereas we are not currently able to provide a proof. And there are examples that cannot be handled by any tool. There is, therefore, room for further improvement. In the near future we plan to add more processors to our current implementation. As remarked above, other processors based on transformations or on reducing the proof obligations in the CTRS problems (for instance, by developing a suitable notion of usable rule or exploiting innermost rewriting [1]) can be added to obtain a more powerful implementation.

References

if-then-else Expressions in the Context-Sensitive Dependency Pair Framework⋆

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Abstract. Syntactic annotations are a powerful mechanism to avoid undesired reductions. In term rewriting, context-sensitive rewriting (CSR) defines the reduction relation which only rewrites the so-called active subterms, i.e., those occurring in replacing arguments of function symbols, as indicated by a replacement map µ, which specifies those arguments. One of the benefits of using CSR is to model the operational behavior of the if-then-else expression in a natural way. Termination of CSR is an interesting problem with several applications in the fields of term rewriting and in the analysis of programming languages like CafeOBJ, Haskell, Maude, OBJ, etc. Direct techniques and frameworks for proving termination of CSR have been developed. But, automatic termination of context-sensitive term rewriting systems (CS-TRSs) involving if-then-else expressions is always hard in practice. In this work, we focus on how if-then-else expressions in CSR are translated into the CS-DP framework, showing why it is difficult to find a termination proof, proposing a new set of context-sensitive usable rules that allows us to simplify termination proofs in these cases and showing some examples where the technique is successfully applied.

Keywords: context-sensitive rewriting, termination, dependency pairs

1 Introduction

Context-Sensitive Rewriting (CSR, [18]) has been shown useful to model evaluation strategies in programming languages. In particular, it is an essential ingredient to analyze the termination behavior of programming languages (like CafeOBJ, Maude, OBJ, etc.), see [14, 20]. In CSR, we start with a pair (R, µ) (often called a CS-TRS) consisting of a term rewriting system (TRS) R and a

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replacement map $\mu$ satisfying $\mu(f) \subseteq \{1, \ldots, \text{ar}(f)\}$ for every function symbol $f$ in the signature $\mathcal{F}$, where $\text{ar}(f)$ is the arity of $f$. Here, $\mu$ is used to discriminate the argument positions on which the rewriting steps are allowed. In this way, we can avoid undesired computations and obtain a terminating behavior for the TRS (with respect to the context-sensitive rewrite relation). One of the most interesting cases where CSR has been shown useful is to model the operational behavior of if-then-else expressions.

Example 1. The following TRS $\mathcal{R}$ [22] provides a definition of factorial which uses CSR for handling the if-then-else expression:

\[
\begin{align*}
0+x & \rightarrow x \quad (1) \\
\text{s}(x)+y & \rightarrow \text{s}(x+y) \quad (2) \\
\text{p}(s(x)) & \rightarrow x \quad (3) \\
\text{if}(\text{true}, x, y) & \rightarrow x \quad (4) \\
\text{if}(\text{false}, x, y) & \rightarrow y \quad (5)
\end{align*}
\]

with $\mu(+) = \mu(\text{s}) = \{1, 2\}$, $\mu(\text{zero}) = \mu(\text{p}) = \mu(\text{if}) = \{1\}$ and $\mu(0) = \mu(\text{true}) = \mu(\text{false}) = \emptyset$. The only function symbol which is restricted by $\mu$ is if, this restriction corresponds to its natural lazy evaluation strategy: the second and third arguments of the expression are not evaluated until the guard (the first argument) is evaluated to true or false. Without the replacement map, this example is non-terminating because $\text{fact}(x)$ calls $\text{fact}(\text{p}(x))$, which then calls $\text{fact}(\text{p}(\text{p}(x)))$ and so on. Thanks to the replacement restrictions, if we want to obtain the factorial of a number $n$ (encoded as $\text{s}^n(0)$):

\[
\text{fact}(\text{s}^n(0)) \rightarrow_\mathcal{R} \text{if}(\text{zero}(\text{s}^n(0)), \text{s}(0), \text{s}^n(0) \ast \text{fact}(\text{p}(\text{s}^n(0)))) \rightarrow_\mathcal{R} \cdots
\]

$\text{zero}(\text{s}^n(0))$ is forced to be reduced first to true or false before the evaluation of the if-then-else expression, thus avoiding undesired reductions until the guard is fully evaluated.

As in unrestricted term rewriting, a notion of context-sensitive dependency pair (CS-DP [2–4]) and several CS-DP frameworks [1, 4, 15] have been defined to capture the possible infinite sequences of calls in a CS-TRS. But, in practice, it is still hard to deal with if-then-else expressions by using existing termination provers. In fact, finding an automatic proof of Example 1, and other examples like [10, Example 1] or [9, Example 3.2.14] are open problems since 1997, 2003 and 2008, respectively.

In this paper, Section 3 analyzes why dealing with if-then-else expressions in CSR is hard in practice, Section 4 imposes the conditions on the notion of chain to allow using the notion of extended basic CS usable rule defined in Section 5, Section 5 proposes a new set of CS usable rules that allows us to simplify termination proofs involving if-then-else expressions and Section 6 shows some examples where the technique is successfully applied.
2 From CSR to the CS-DP Framework

See [8] and [18] for basics on term rewriting and CSR, respectively. Throughout the paper, $\mathcal{X}$ denotes a countable set of variables and $\mathcal{F}$ denotes a signature, i.e., a set of function symbols each having a fixed arity given by a mapping $ar : \mathcal{F} \to \mathbb{N}$. The set of terms built from $\mathcal{F}$ and $\mathcal{X}$ is $\mathcal{T}(\mathcal{F}, \mathcal{X})$. Terms are viewed as labelled trees in the usual way. The symbol labeling the root of the term $s$ is denoted as $root(s)$. Positions $p, q, \ldots$ are represented by chains of positive natural numbers used to address subterms of $s$. Given positions $p, q$, we denote its concatenation as $p.q$. We denote the empty chain by $\varepsilon$. The set of positions of a term $s$ is $\mathcal{P}(s)$.

For a replacement map $\mu$, the set of active positions $\mathcal{P}(s)^\mu(s)$ of $s \in \mathcal{T}(\mathcal{F}, \mathcal{X})$ is:

$$\mathcal{P}(s)^\mu(s) = \{\varepsilon\} \cup \bigcup_{i \in \mathbb{N}} i.\mathcal{P}(s)\{\varepsilon\}, \text{if } s \notin \mathcal{X}.$$  

We say that $s \mu \triangleright t$ if there is $p \in \mathcal{P}(s)^\mu(s)$ such that $t = s|_p$. We write $s \triangleright \mu \mu t$ if $s \mu \triangleright t$ and $s \neq t$. Moreover, $s \triangleright \mu \mu t$ if there is a frozen position $p$, i.e. $p \in \mathcal{P}(s)\mathcal{T}(\mathcal{F}, \mathcal{X})$ such that $t = s|_p$. Let $\text{Var}(s) = \{x \in \mathcal{X} \mid \exists p \in \mathcal{P}(s)\mathcal{T}(\mathcal{F}, \mathcal{X}), \text{Var}^\mu(s) = \{x \in \text{Var}(s) \mid \exists p \in \mathcal{P}(s)\mathcal{T}(\mathcal{F}, \mathcal{X}), s|_p = x\}$ and $\text{Var}^\mu(s) = \{x \in \text{Var}(s) \mid s \triangleright \mu \mu t\}$. A rewrite rule is an ordered pair $(\ell, r)$, written $\ell \rightarrow r$, with $\ell, r \in \mathcal{T}(\mathcal{F}, \mathcal{X})$, $\ell \notin \mathcal{X}$ and $\text{Var}(r) \subseteq \text{Var}(\ell)$. A TRS is a pair $\mathcal{R} = (\mathcal{F}, \mathcal{R})$ where $\mathcal{R}$ is a set of rewrite rules. Given $\mathcal{R} = (\mathcal{F}, \mathcal{R})$, we consider $\mathcal{F}$ as the disjoint union $\mathcal{F} = \mathcal{C} \uplus \mathcal{D}$ of symbols $c \in \mathcal{C}$, called constructors and symbols $f \in \mathcal{D}$, called defined functions, where $\mathcal{D} = \{\text{root}(\ell) \mid \ell \rightarrow r \in \mathcal{R}\}$ and $\mathcal{C} = \mathcal{F} - \mathcal{D}$. Given a CS-TRS $(\mathcal{R}, \mu)$, we have $s \mu \rightarrow_{\mathcal{R}, \mu} t$ (alternative $s \mu \rightarrow_{\mathcal{R}, \mu} t$ if we want to make the position explicit) if there are $\ell \rightarrow r \in \mathcal{R}$, $p \in \mathcal{P}(s)^\mu(s)$ and a substitution $\sigma$ with $s|_p = \sigma(\ell)$ and $t = s|_{\sigma(r)}|_p$. A CS-TRS $(\mathcal{R}, \mu)$ is terminating if $\rightarrow_{\mathcal{R}, \mu}$ is well-founded.

Dependency pairs [7] describe the propagation of function calls in rewrite sequences. The main idea for using dependency pairs when modeling infinite rewrite sequences issued from (finite) TRSs is simple: for every infinite rewrite sequence there is a finite set of function calls (often called a cycle) which are used infinitely often. In CSR, we have two kinds of (potential) function calls: direct calls, i.e., calls at active (replacing) positions and delayed calls, i.e., calls at frozen (non-replacing) positions that can be activated in forthcoming reduction steps. These function calls are captured in two different ways. For rules $\ell \rightarrow r$ such that $r$ contains an active subterm $v = g(s_1, \ldots, s_m)$ where $g$ is a defined function, the function call to $g$ is represented as a new rule $u \rightarrow v$ (called dependency pair) where $u = f(\ell_1, \ldots, \ell_k)$ if $\ell = f(\ell_1, \ldots, \ell_k)$ and $v = g^2(s_1, \ldots, s_m)$. The notation $f^2$ for a given symbol $f$ means that $f$ is marked. In practice, we often capitalize $f$ and use $F$ instead of $f^2$ in our examples. Function calls to $g$ which are at frozen positions of $r$ cannot be issued ‘immediately’, but could be activated ‘in the future’. This situation is carried out by the migrating variables and modeled by the collapsing DPs. Given a rule $\ell \rightarrow r$, $x$ is a migrating variable if $x$ is at an active position in $r$ but not in $\ell$ [2]. For rules $\ell \rightarrow r$, collapsing DPs are pairs of the form $u \rightarrow x$ where $u = f(\ell_1, \ldots, \ell_k)$ if $\ell = f(\ell_1, \ldots, \ell_k)$ and $x$ is a migrating variable. The idea is that calls which can eventually be activated are subterms of $\sigma(x)$ for $\sigma$ being the matching substitution of the rewriting step involving the rule $\ell \rightarrow r$. Furthermore, $\sigma(x)$ is of the form $\theta(C[s])$, being $C$ a
hiding context, s a hidden term and \( \theta \) a substitution. We say that \( s \) is a hidden term if root\((s) \in D \) and there exists a rule \( \ell \to r \) with \( r \vdash \mu s \). Given a hidden term of the form \( g(s_1, \ldots, s_m) \), a function call activation is represented by a rule \( g(s_1, \ldots, s_m) \to g^\sharp(s_1, \ldots, s_m) \) in the unhiding TRS, where the marking represents the activation. A function symbol \( f \) hides position \( i \) in the rule \( \ell \to r \) if \( r \vdash \mu f(r_1, \ldots, r_n) \) and there is \( i \in \mu(f) \) and a term \( s \) such that \( r_i \vdash \mu s \). If root\((s) \in D \) or \( s \) is a variable in \( (\text{Var}^\sharp(\ell) \cap \text{Var}^\sharp(r)) - (\text{Var}^\mu(\ell) \cup \text{Var}^\mu(r)) \). A context \( C[\square] \) is hiding if \( C[\square] = \square \), or \( C[\square] = f(t_1, \ldots, t_{i-1}, C'[\square], t_{i+1}, \ldots, t_k) \), where \( f \) hides position \( i \) and \( C'[\square] \) is a hiding context. If \( f \) hides position \( i \), the unhiding TRS contains a rule of the form \( f(x_1, \ldots, x_{i-1}, x_i, x_{i+1}, \ldots, x_k) \to x_i \) to extract any possible hiding context from a term. Formally, \( \text{DP}(R, \mu) = \text{DP}_F(R, \mu) \cup \text{DP}_X(R, \mu) \) where \( \text{DP}_F(R, \mu) = \{ f^\sharp(t_1, \ldots, t_k) \to g^\sharp(s_1, \ldots, s_m) \mid f(t_1, \ldots, t_k) \to r \in R, r \vdash \mu g(s_1, \ldots, s_m), g \in D, f(t_1, \ldots, t_k) \vdash \mu g(s_1, \ldots, s_m) \} \).

To prove termination, we have to show that there is no infinite chain of CS-DPs. A sequence \( u_1 \to v_1, u_2 \to v_2, \ldots \) of CS-DPs is a chain if there is a substitution \( \sigma \) such that for all \( i \geq 1 \), (1) if \( u_i \to v_i \) is not collapsing, then \( \sigma(u_i) \to_{\mu, R, \mu} \sigma(u_{i+1}) \) or (2) if \( u_i \to v_i \) is collapsing, then there is a term \( w_i \) such that (2a) \( \sigma(u_i) \stackrel{\circ}{\to}_{\text{unhide}(R, \mu)} \sigma(w_i) \) and (2b) \( \sigma(v_i) \stackrel{\circ}{\to}_{\text{unhide}(R, \mu)} \sigma(u_{i+1}) \).

**Example 2.** For \((R, \mu)\) in Example 1, we obtain the following CS-DPs:

\[
\begin{align*}
s(x)+^\sharp y & \to x+^\sharp y \quad (12) & \text{FACT}(x) & \to \text{ZERO}(x) \quad (16) \\
s(x)*^\sharp y & \to y+^\sharp (x*y) \quad (13) & \text{IF}(\text{true}, x, y) & \to x \quad (17) \\
s(x)*^\sharp y & \to x*y \quad (14) & \text{IF}(\text{false}, x, y) & \to y \quad (18) \\
\text{FACT}(x) & \to \text{IF}(\text{zero}(x), s(0), x*\text{fact}(p(x)))) \quad (15)
\end{align*}
\]

with \( \mu \) extended by \( \mu(+) = \mu(\cdot) = \{1, 2\} \), and \( \mu(\text{ZERO}) = \mu(\text{FACT}) = \mu(\text{IF}) = \{1\} \). DPs (12)-(16) capture the direct calls and collapsing DPs (17)-(18) capture the delayed calls. The unhiding TRS contains the following rules:

\[
\begin{align*}
x*y & \to y \quad (19) & \text{fact}(p(x)) & \to \text{FACT}(p(x)) \quad (20)
\end{align*}
\]

obtained from the rule (8), where \( * \) hides position 2 and \( \text{fact}(p(x)) \) is a hidden term. An example of a CS- DP chain involving the collapsing pair (18) is:

\[
\begin{align*}
\text{FACT}(\text{s}(x)) & \to_{(15)} \text{IF}(\text{zero(s}(x)), s(0), s(x)*\text{fact}(p(s(x)))) \to_{(1)} \\
\text{IF}(\text{false}, s(0), s(x)*\text{fact}(p(s(x)))) & \to_{(18)} s(x)*\text{fact}(p(s(x))) \to_{(19)} \text{fact}(p(s(x))) \to_{(20)} \\
\text{FACT}(p(s(x))) & \to_{(1)} \ldots
\end{align*}
\]

where the rule (19) removes the hiding context and the rule (20) activates the call to continue the chain.
Termination is characterized by the absence of infinite chains [2, 4]. The CS-DP framework helps us to mechanize the process of proving the absence of minimal chains. As usual in the DP framework, we abstract the previous characterization of termination of CS-TRS for a CS-TRS \((R, \mu)\) with the set of CS-DPs \(DP(R, \mu)\) using generic TRSs \(P\), \(R\) and \(S\).

**Definition 1 (Chain of Pairs [15]).** Let \(R\), \(P\) and \(S\) be TRSs and \(\mu\) a replacement map where \(S = S_{\varphi_{\mu}} \sqcup S_{\varphi}\), \(S_{\varphi_{\mu}}\) are rules of the form \(s \rightarrow t \in S\) such that \(s \varphi_{\mu} t\) and \(S_{\varphi} = S - S_{\varphi_{\mu}}\). A \((P, R, S, \mu)\)-chain is a finite or infinite sequence of pairs \(u_i \rightarrow v_i \in P\), together with a substitution \(\sigma\) satisfying that, for all \(i \geq 1\),

1. if \(v_i \notin \text{Var}(u_i) - \text{Var}^\mu(u_i)\), then \(\sigma(v_i) = w_i \varphi^*_{R, \mu} \sigma(u_{i+1})\), and
2. if \(v_i \in \text{Var}(u_i) - \text{Var}^\mu(u_i)\), then \(\sigma(v_i) \varphi^*_{S_{\varphi_{\mu}}} \circ \varphi_{S_{\varphi}} w_i \varphi^*_{R, \mu} \sigma(u_{i+1})\).

An infinite \((P, R, S, \mu)\)-chain is called minimal if for all \(i \geq 1\), \(w_i\) is \((R, \mu)\)-terminating.

In Definition 1, \(P\) plays the role of \(DP(R, \mu)\) and \(S\) has two components \(S_{\varphi_{\mu}}\) and \(S_{\varphi}\) which are useful to model \((2a)\) above when connecting a collapsing pair to another pair.

To use the CS-DP framework to prove the absence of minimal chains of CS-DPs, we let \(P\) be the set of CS-DPs, \(R\) the initial TRS, and \(S\) the unhiding TRS.

A CS problem \((P, R, S, \mu)\) is finite if there is no infinite minimal \((P, R, S, \mu)\)-chain. A CS processor \(\text{Proc}\) is a mapping from CS problems into sets of CS problems. A CS-processor \(\text{Proc}\) is sound if for all CS problems \(\tau\), \(\tau\) is finite whenever \(\forall \tau' \in \text{Proc}(\tau)\), \(\tau'\) is finite.\(^1\).

**Theorem 1 (CS-DP Framework [15]).** Let \((R, \mu)\) be a CS-TRS. We construct a tree whose nodes are labeled with CS problems or “yes”, and whose root is labeled with \((DP(R, \mu), R, \text{unh}(R, \mu), \mu^1)\). For every inner node labeled with \(\tau\), there is a sound processor \(\text{Proc}\) satisfying one of the following conditions:

1. \(\text{Proc}(\tau) = \emptyset\) and the node has just one child, labeled with “yes”.
2. \(\text{Proc}(\tau) \neq \emptyset\), and the children of the node are labeled with the CS problems in \(\text{Proc}(\tau)\).

If all leaves of the tree are labeled with “yes”, then \((R, \mu)\) is terminating.

From now on, we assume that in all the considered CS problems \((P, R, S, \mu)\) the following property holds: for all \(s \rightarrow t \in S\) and substitution \(\sigma\), if \(\sigma(t)\) generates an infinite minimal \((P, R, S, \mu)\)-chain then \(\sigma(s)\) is not \((R, \mu)\)-terminating. This property always holds on the root CS problem created in Theorem 1.

\(^1\) In order to keep our presentation simpler, we do not introduce here the notions related with completeness of processors, needed for non-termination proofs.
3 Analyzing if-then-else Expressions

One of the most powerful processors to deal with CS problems is the μ-reduction pair processor, a processor that discards pairs that can be strictly oriented using orderings. A μ-reduction pair \((\succeq, \sqsubset)\) consists of a stable and μ-monotonic\(^2\) quasi-ordering \(\succeq\), and a well-founded stable relation \(\sqsubset\) on terms in \(T(F, X)\) which are compatible, i.e., \(\succeq \circ \sqsubset\) or \(\sqsubset \circ \succeq\) [2]. Given a CS problem \(\tau = (P, R, S, \mu)\), if there is a μ-reduction pair such that \(P \cup S \subseteq \epsilon \sqsubset \sqsubset\) and \(R \subseteq \succeq\) then \((P, R, S, \mu)\) is finite if \((P - P_{\sqsubset}, R, S - S_{\sqsubset}, \mu)\) is finite, where \(P_{\sqsubset}\) and \(S_{\sqsubset}\) represent the set of rules from \(P\) and \(S\) oriented using \(\sqsubset\). The μ-reduction pair processor can be improved using the notion of usable rule [6]. Usable rules, initially connected to innermost termination, allow us to discard those rules from \(R\) that are not directly involved in (possible) infinite minimal \((P, R, S, \mu)\)-chains. In rewriting (and also in CSR), the notion of usable rule is connected with \(\mathcal{C}_s\)-termination [12, 21]. A TRS \(\mathcal{R} = (F, R)\) is \(\mathcal{C}_s\)-terminating if \(R \triangleright \mathcal{C}_s\) is terminating, where \(\mathcal{C}_s = \{c(x, y) \rightarrow x, c(x, y) \rightarrow y\} \) (with \(c \notin F\)). The idea behind the usable rules is that for every infinite \((P, R, S, \mu)\)-chain we can construct an infinite sequence where rewriting steps using \(\mathcal{R}\) can be simulated by rewriting steps using \(\mathcal{U}_s(R)\) and \(\mathcal{C}_s\), where \(\mathcal{U}_s(R)\) is the set of usable rules of \(\tau\). So, instead of \(\mathcal{R} \subseteq \mathcal{C}_s\), we only need \(\mathcal{U}_s(R) \triangleright \mathcal{C}_s\) and \(\mathcal{C}_s\).

In [16], the notion of CS usable rule was given for chains of pairs. This notion is different from the one given in standard rewriting. For example, if we consider the following CS problem \(\tau_1 = \{\{15\}, \{17\}, \{18\}\}, \mathcal{R}, \{(19), (20)\}, \mu\}\) obtained from Example 2, the set of CS usable rules in \(\tau_1\) is \(\mathcal{R}\). This is caused by the presence of migrating variables. In the presence of migrating variables, every rule headed by a symbol appeared at a frozen positions in the right-hand side of a rule in \(\mathcal{R}\) is considered usable (in this case *, fact and p, and by transitivity +, if and zero).

But, if we look closely at the CS-DP chain \(\text{FACT}(s(x)) \rightarrow_{15} \cdots \rightarrow_{20} \text{FACT}(p(s(x))) \rightarrow_{\mathcal{R}, \mu} \cdots\) from Example 2, we notice that \(x\) appears at an active position in \(\text{FACT}(s(x))\) and \(\text{FACT}(p(s(x)))\), i.e., the variable \(x\) does not behave in the CS-DP chain as a migrating variable. This “conservative” behavior is shared by all the if-then-else unresolved expressions in the literature.

In [16], a notion of CS usable rule that allows us to obtain a smaller (and closer to the unrestricted rewriting approach) set of CS usable rules was also given, but it is restricted to pairs and rules that are strongly conservative. A rule \(t \rightarrow r\) is strongly conservative if it is conservative (i.e., \(\text{Var}^+(r) \subseteq \text{Var}^+(t))\), and \(\text{Var}^+(t) \cap \text{Var}^+(\ell) = \text{Var}^+(r) \cap \text{Var}^+(\ell) = \emptyset\). This notion is hard in practice and discards, by definition, CS problems with collapsing pairs. So, we have to relax strongly conservative conditions in the presence of collapsing pairs. We carry out this task in two steps:

\(^2\) A binary relation \(R\) on terms is μ-monotonic if for all terms \(s, t, t_1, \ldots, t_m\), and \(m\)-ary symbols \(f\), whenever \(s R t\) and \(i \in \mu(f)\) we have \(f(t_1, \ldots, t_{i-1}, s, \ldots, t_m) R f(t_1, \ldots, t_{i-1}, t, \ldots, t_m)\).
1. We introduce the notion of strongly minimal \((\mathcal{P}, \mathcal{R}, \mathcal{S}, \mu)\)-chain. This notion allows us to ensure that possible infinite delayed calls are always introduced by means of pairs or rules, i.e., the initial term of an infinite \((\mathcal{P}, \mathcal{R}, \mathcal{S}, \mu)\)-chain does not contain subterms that can generate infinite \((\mathcal{P}, \mathcal{R}, \mathcal{S}, \mu)\)-chains.

2. We find conditions to obtain a conservative behavior in the presence of collapsing pairs to reduce the set of CS usable rules.

## 4 Strongly Minimal \((\mathcal{P}, \mathcal{R}, \mathcal{S}, \mu)\)-Chains

From a termination point of view, terms can be divided naturally into two sets: terminating and non-terminating terms. Among non-terminating terms we can extract those that are minimal, namely, those that are non-terminating and do not contain non-terminating proper subterms. By construction, minimal non-terminating terms are headed by a defined symbol, i.e., they are function calls. For each term in an infinite rewrite sequence we can always extract a subterm that is minimal non-terminating. The notion of dependency pair \([7]\) captures the shape of those minimal non-terminating terms that are introduced in the sequence by means of rewriting rules: instantiated function calls that appear on the right-hand side of the rules.

In CSR, we use two notions of minimality \([4]\):

- a term is strongly minimal non-\(\mu\)-terminating, \(t \in T_{\infty, \mu}\) if it is non-\(\mu\)-terminating and does not contain non-\(\mu\)-terminating proper subterms.
- a term is minimal non-\(\mu\)-terminating, \(t \in M_{\infty, \mu}\) if it is non-\(\mu\)-terminating and does not contain proper subterms at active positions that are non-\(\mu\)-terminating.

The second notion is used in the CS-DP framework to formalize the notion of minimal \((\mathcal{P}, \mathcal{R}, \mathcal{S}, \mu)\)-chain. Infinite minimal \(\mu\)-rewrite sequences starting by strongly minimal non-\(\mu\)-terminating terms are paramount in the latest improvements of the CS-DP framework (mechanization of the unhiding pairs \([1, 4, 15]\)) and allow us to discard many problematic infinite minimal \(\mu\)-rewrite sequences, but until now the CS-DP framework does not take advantage of this notion. We introduce this notion in the following definition.

**Definition 2.** An infinite \((\mathcal{P}, \mathcal{R}, \mathcal{S}, \mu)\)-chain \(u_1 \rightarrow v_1, u_2 \rightarrow v_2, \ldots\) is called strongly minimal if it is minimal and there is no rule \(s \rightarrow t \in \mathcal{S}\) and substitutions \(\sigma, \theta\) such that \(\sigma(u_1) \triangleright \theta(s)\) and \(\theta(t)\) starts an infinite minimal \((\mathcal{P}, \mathcal{R}, \mathcal{S}, \mu)\)-chain.

The following result allows us to use strongly minimal chains in the CS-DP framework.

**Theorem 2.** Let \(\tau = (\mathcal{P}, \mathcal{R}, \mathcal{S}, \mu)\) be a CS problem such that for every \(s \rightarrow t \in \mathcal{S}, s = f(s_1, \ldots, s_m)\) and \(t = g(s_1, \ldots, s_m)\). Then, \(\tau\) is finite if there is no infinite strongly minimal \((\mathcal{P}, \mathcal{R}, \mathcal{S}, \mu)\)-chain.
Proof. If $\tau$ is not finite, then there is an infinite minimal $\langle P, R, S, \mu \rangle$-chain of pairs $u_1 \rightarrow v_1, u_2 \rightarrow v_2, \ldots$ together with a substitution $\sigma$. By structural induction on $\sigma(u_1)$:

1. If there is no $u'_1$, rule $s \rightarrow t \in S_t$ and substitution $\theta$ such that $\sigma(u_1) \triangleright u'_1$ and $u'_1 = \theta(s)$, then the infinite minimal $\langle P, R, S, \mu \rangle$-chain is strongly minimal.

2. If there is $u'_1$, rule $s \rightarrow t \in S_t$ and substitution $\theta$ such that $\sigma(u_1) \triangleright u'_1$ and $u'_1 = \theta(s)$, then without loss of generality we can choose $u'_1$ to be minimal (i.e., $u'_1$ has no subterm satisfying the previous conditions) and there is an infinite minimal $\langle P, R, S, \mu \rangle$-chain starting from $\theta(t)$, by Definition 2. The obtained infinite minimal $\langle P, R, S, \mu \rangle$-chain is strongly minimal, if not there is $u''_1$ such that $\theta(t) \triangleright u''_1$. But, by hypothesis, $\theta(s) \triangleright u''_1$ contradicting the minimality of $u'_1$.

Rules appeared in $\text{unh}_2(R, \mu)$ always satisfy the condition imposed on $S_t$ in Theorem 2.

5 Relaxing CS Usable Rules Constraints

The goal in this section is to obtain a smaller set of CS usable rules in the presence of collapsing pairs. We will do it by means of conservativity and left-linearity conditions. Left linearity together with conservativity is useful in connection with CS usable rules because ensures that a term occurring at a frozen position does not show up at an active position by means of a variable instantiation after pair or rule applications. Collapsing pairs are not conservative, but if we have an infinite minimal $\langle P, R, S, \mu \rangle$-chain with a collapsing pair $u_i \rightarrow v_i$, we know that $\sigma(v_i) \xrightarrow{\sigma(S_t)} \circ \xrightarrow{\sigma(S_t)} t_i \rightarrow \tau^\mu \circ \epsilon \rightarrow S \circ t_i \rightarrow \tau^\mu \circ \sigma(u_{i+1})$. So, if we ensure that any instantiated term occurring at an active position in the right-hand side of pairs $S_t$ comes from a term that (previously in the chain) was at an active position, we say now that the CS problem is conservative with respect to $S$.

**Definition 3 (Conditions for $S$).** Let $\tau = \langle P, R, S, \mu \rangle$ be a CS problem. We say that $\tau$ is conservative with respect to $S$ if $S$ is conservative and the following conditions hold:

- for all $s \rightarrow t \in S_t$, $s = f(s_1, \ldots, s_m)$ and $t = g(s_1, \ldots, s_m)$,
- for each $s \rightarrow t \in S_t$ and for each $u \rightarrow v \in P \cup R$, if there is a non-variable subterm $v'$ of $v$ at a frozen position such that $\theta = \text{mgut}(v', s)$, then $v' = s$ up to renaming of variables and $u \rightarrow v'$ must be conservative.

These conditions force that all the delayed calls are hidden terms in $P$ or $R$. Our goal is to impose conservativity and left-linearity instead of strong conservativity, thus relaxing the requirements on the right-hand sides of the rules. The payoff is that we still have to add as usable those rules from symbols that occur at frozen positions in the right-hand side of rules and pairs, as we can see in the counterexample from [16, Example 4].
Example 3. Consider the following CS problem \( \tau = (P, R, S, \mu) \) taken from [16, Example 4] that satisfies the conservativity and left-linearity conditions:

\[
P = \{ A(x, y) \rightarrow B(x, x), \ B(x, c) \rightarrow D(x, x), \ D(x, e) \rightarrow A(x, x) \}
\]

\[
R = \{ a(x, y) \rightarrow b(x, x), \ b(x, c) \rightarrow d(x, x), \ d(x, e) \rightarrow a(x, x), \ c \rightarrow e \}
\]

and \( S = \emptyset \), with \( \mu(A) = \mu(a) = \mu(B) = \mu(b) = \{1, 2\} \) and \( \mu(c) = \mu(e) = \emptyset \). If we do not consider symbols occurring at frozen positions in the left-hand side of pairs and rules, we have no usable rules. But, the following infinite \((P, R, S, \mu)\)-chain exists and the rule \( c \rightarrow e \) is usable

\[
A(c, c) \leftarrow P, \ B(c, c) \leftarrow P, \ D(c, e) \leftarrow R, \ D(c, e) \leftarrow R, \ A(c, c) \leftarrow \cdots
\]

But, as remarked in [16], this is not a hard restriction in practice because most programs are constructor systems, that is, no defined symbols occur below the root in the left-hand side of the rules. Therefore, we define the extended basic \( \mu \)-dependency considering this extension.

Let \( \text{Fun}^\mu(s) \) be the set of \( \mu \)-replacing symbols in a term \( s \in T(F, X) \), \( \text{Fun}^\mu(s) = \{ f \mid \exists p \in \text{Pos}^\mu(s), f = \text{root}(s|_p) \} \), and \( \text{Fun}^\emptyset(s) \) the set of non-\( \mu \)-replacing symbols in a term \( s \in T(F, X) \), \( \text{Fun}^\emptyset(s) = \{ f \mid \exists p \in \text{Pos}(s) - \text{Pos}^\mu(s), f = \text{root}(s|_p) \} \). Let \( \text{Rls}_\mu(f) = \{ \ell \rightarrow r \in R \mid \text{root}(\ell) = f \} \).

**Definition 4 (Extended Basic \( \mu \)-Dependency).** Given a TRS \((F, R)\) and a replacement map \( \mu \), we say that \( f \in F \) has an extended basic \( \mu \)-dependency on \( h \in F \), written \( f \triangleright_{R, \mu} h \), if \( f = h \) or there is a function symbol \( g \) with \( g \triangleright_{R, \mu} h \) and a rule \( \ell \rightarrow r \in \text{Rls}_\mu(f) \) with \( g \in \text{Fun}^\emptyset(\ell) \cup \text{Fun}^\mu(\ell) \) and \( r \)

We extend now the notion of usable rule accordingly.

**Definition 5 (Extended Basic CS Usable Rules).** Let \( \tau = (P, R, S, \mu) \) be a CS problem. The set \( U^\mu_\tau(R) \) of extended basic context-sensitive usable rules of \( \tau \) is

\[
U^\mu_\tau(R) = \bigcup_{u \rightarrow v \in \mathcal{P} \cup \mathcal{S}, f \in \text{Fun}^\mu(u) \cup \text{Fun}^\emptyset(v), f \triangleright_{R, \mu} g} \text{Rls}_\mu(g)
\]

We obtain the processor \( \text{Proc}_{UR} \). The pairs \( P \) in a CS problem \((P, R, S, \mu)\), where \( P \) is a TRS over the signature \( \mathcal{G} \), are partitioned as follows: \( \mathcal{P}_X = \{ u \rightarrow v \in \mathcal{P} \mid v \in \text{Var}(u) \sim \text{Var}^\mu(u) \} \) and \( \mathcal{P}_G = \mathcal{P} - \mathcal{P}_X \).

**Theorem 3.** Let \( \tau = (P, R, S, \mu) \) be a CS problem such that \((a)\) \( P_X \cup U^\mu_\tau(R) \cup S_{\tau, \mu} \) is conservative and left-linear, and \((b)\) whenever \( P_X \neq \emptyset \) we have that \( P_X \) is left-linear and \( \tau \) is conservative with respect to \( S \). Let \( (\mathcal{S} \supseteq \mathcal{S}) \) be a \( \mu \)-reduction pair such that

1. \( P \subseteq \mathcal{S} \cup \mathcal{R} \cup U^\mu_\tau(R) \subseteq \mathcal{S} \),
2. whenever \( P_X \neq \emptyset \) we have that \( S \subseteq \mathcal{S} \cup \mathcal{R} \).

Let \( P_{\subseteq} = \{ u \rightarrow v \in \mathcal{P} \mid u \sqsubseteq v \} \) and \( S_{\subseteq} = \{ s \rightarrow t \in \mathcal{S} \mid s \sqsubseteq t \} \). Then, the processor \( \text{Proc}_{UR} \) given by

\[
\text{Proc}_{UR}(\tau) = \begin{cases} \{(P - P_{\subseteq}, R, S - S_{\subseteq}, \mu)\} & \text{if (1) and (2) hold} \\ \{(P, R, S, \mu)\} & \text{otherwise} \end{cases}
\]
is sound.

The proof is at the end of the section, we introduce now the partial results needed to obtain it. We use the interpretation given in [16, Definition 11] to define the new interpretation.

**Definition 6 (Basic \(\mu\)-Interpretation [16]).** Let \((R, \mu)\) be a CS-TRS over the signature \(\mathcal{F}\) and \(\Delta \subseteq \mathcal{F}\). Let \(>\) be an arbitrary total ordering on terms in \(\mathcal{T}(\mathcal{F} \cup \{\perp, \cdot\}, \mathcal{X})\) where \(\perp\) is a fresh constant symbol and \(\cdot\) is a fresh binary symbol. The basic \(\mu\)-interpretation \(I_{0, \Delta, \mu}\) is a mapping from \(\mu\)-terminating terms in \(\mathcal{T}(\mathcal{F}, \mathcal{X})\) to terms in \(\mathcal{T}(\mathcal{F} \cup \{\perp, \cdot\}, \mathcal{X})\) defined as follows:

\[
I_{0, \Delta, \mu}(t) = \begin{cases} 
  t & \text{if } t \in \mathcal{X} \\
  f(I_{0, \Delta, \mu}(t_1), \ldots, I_{0, \Delta, \mu}(t_n)) & \text{if } t = f(t_1, \ldots, t_n) \text{ and } f \in \Delta \\
  c(f(I_{0, \Delta, \mu}(t_1), \ldots, I_{0, \Delta, \mu}(t_k)), t') & \text{if } t = f(t_1, \ldots, t_k) \text{ and } f \notin \Delta 
\end{cases}
\]

where \(I_{0, \Delta, \mu, t}(t) = \begin{cases} 
  I_{0, \Delta, \mu}(t) & \text{if } i \in \mu(f) \\
  t & \text{if } i \notin \mu(f)
\end{cases}\)

\(t' = \text{order }\{\{I_{0, \Delta, \mu}(u) \mid t \mapsto_{\Delta, \mu} u\}\}\)

\(\text{order}(T) = \begin{cases} 
  \perp, & \text{if } T = \emptyset \\
  c(t, \text{order}(T - \{t\})) & \text{if } t \text{ is minimal in } T \text{ w.r.t. } >
\end{cases}\)

**Lemma 1.** [16] Let \((R, \mu)\) be a CS-TRS over the signature \(\mathcal{F}\) and \(t\) in \(\mathcal{T}(\mathcal{F}, \mathcal{X})\). If \(t\) is \(\mu\)-terminating then \(I_{0, \Delta, \mu}\) is well-defined.

But, in order to deal with collapsing pairs, we allow that at frozen positions we can have terms that are interpreted (or partially interpreted). Then, to have a unique interpretation for each term we have to parametrize it with respect to an infinite strongly minimal \((P, R, S, \mu)\)-chain.

**Definition 7 (Extended Basic \(\mu\)-Interpretation).** Let \((P, R, S, \mu)\) be a CS problem where \(R\) is a TRS over the signature \(\mathcal{F}\) and \(\Delta \subseteq \mathcal{F}\). Let \(A\) be an infinite strongly minimal \((P, R, S, \mu)\)-chain of the form \(u_1 \rightarrow v_1, u_2 \rightarrow v_2, \ldots\) Let \(I_{\Delta, \mu, A}\) be an interpretation that satisfies:

- \(I_{\Delta, \mu, A}(\sigma(u_1)) = I_{0, \Delta, \mu}(\sigma(u_1))\), and
- if there is a pair \(u_i \rightarrow v_i \in P\) and a variable \(x \in \text{Var}(u_i) \cap \text{Var}(v_i)\), \(I_{\Delta, \mu, A}(\sigma(x))\) has the same interpretation in \(u_i\) and \(v_i\).
- if there is a rule \(\ell_i \rightarrow r_i \in R\) and a variable \(x \in \text{Var}(\ell_i) \cap \text{Var}(r_i)\), \(I_{\Delta, \mu, A}(\sigma(x))\) has the same interpretation in \(\ell_i\) and \(r_i\).

**Definition 8.** Let \((P, R, S, \mu)\) be a CS problem where \(R\) is a TRS over the signature \(\mathcal{F}\) and \(\Delta \subseteq \mathcal{F}\). Let \(A\) be an infinite strongly minimal \((P, R, S, \mu)\)-chain. We denote by \(\sigma_{I_{\Delta, \mu, A}}\) a substitution that replaces occurrences of \(x \in \text{Var}(t)\) by \(I_{\Delta, \mu, A}(\sigma(x))\).
Lemma 2. Let \((P, R, S, \mu)\) be a CS problem where \(R\) is a TRS over the signature \(F\) and \(\Delta \subseteq F\). Let \(A\) be an infinite strongly minimal \((P, R, S, \mu)\)-chain. Let \(t\) be a linear term and \(\sigma\) be a substitution. If all subterms \(t'\) of \(t\) at frozen positions are from \(\mathcal{T}(\Delta, X)\) and \(\sigma(t)\) is \((R, \mu)\)-terminating, then \(\mathcal{I}_{\Delta, \mu, A}(\sigma(t)) \Rightarrow_{\mathcal{C}_t, \mu} \mathcal{I}_{\Delta, \mu, A}(t)\). If \(t\) only contain \(\Delta\)-symbols at active positions, then we have \(\mathcal{I}_{\Delta, \mu, A}(\sigma(t)) = \mathcal{I}_{\Delta, \mu, A}(t)\).

Proof. By structural induction on \(t\):

- If \(t\) is a variable then \(\mathcal{I}_{\Delta, \mu, A}(\sigma(t)) = \sigma_{\mathcal{I}_{\Delta, \mu, A}(t)}\).
- If \(t = f(t_1', \ldots, t_k')\) then
  - If \(f \in \Delta\) then \(\mathcal{I}_{\Delta, \mu, A}(\sigma(t)) = f(\sigma(t_1'), \ldots, \sigma(t_k'))\). Terms \(\sigma(t_i')\) are \(\mu\)-terminating for \(i \in \mu\)(f). By induction hypothesis, for all terms \(t_i'\) s.t. \(i \in \mu\)(f), we have \(\mathcal{I}_{\Delta, \mu, A}(\sigma(t_i')) \Rightarrow_{\mathcal{C}_t, \mu} \sigma_{\mathcal{I}_{\Delta, \mu, A}(t_i')}\). And for all \(t_i'\) s.t. \(i \notin \mu\)(f), we have that \(t_i'\) only contains \(\Delta\) symbols, then \(t_i' = \mathcal{I}_{\Delta, \mu, A}(\sigma(t_i')) = \sigma_{\mathcal{I}_{\Delta, \mu, A}(t_i')}\). This implies \(f(\sigma(t_1'), \ldots, \sigma(t_k')) \Rightarrow_{\mathcal{C}_t, \mu} \sigma_{\mathcal{I}_{\Delta, \mu, A}(t)}\).
  - If \(f \notin \Delta, \mathcal{I}_{\Delta, \mu, A}(\sigma(t)) \Rightarrow_{\mathcal{C}_t, \mu} \tau(f(t_1', \ldots, t_k'), t')\) for some \(t'\). Applying a \(\mathcal{C}_t\) step to this term, we obtain again the term \(f(t_1', \ldots, t_k')\), and using the previous item result, we get \(f(t_1', \ldots, t_k') \Rightarrow_{\mathcal{C}_t, \mu} \sigma_{\mathcal{I}_{\Delta, \mu, A}(t)}\).

Then we conclude \(\mathcal{I}_{\Delta, \mu, A}(\sigma(t)) \Rightarrow_{\mathcal{C}_t, \mu} \sigma_{\mathcal{I}_{\Delta, \mu, A}(t)}\).

The second part of the lemma is proved similarly. If \(t\) is a variable then \(\mathcal{I}_{\Delta, \mu, A}(\sigma(t)) = \sigma_{\mathcal{I}_{\Delta, \mu, A}(t)}\). Now let \(t = f(t_1, \ldots, t_k)\). Since \(f \in \Delta, \mathcal{I}_{\Delta, \mu, A}(\sigma(t)) = \mathcal{I}_{\Delta, \mu, A}(f(\sigma(t_1), \ldots, \sigma(t_k))) = f(\sigma(t_1'), \ldots, \sigma(t_k'))\). For \(i \in \mu\)(f), we have \(\mathcal{I}_{\Delta, \mu, A}(\sigma(t_i')) = \sigma_{\mathcal{I}_{\Delta, \mu, A}(t_i')}\) by the induction hypothesis. For \(i \notin \mu\)(f), we have \(\mathcal{I}_{\Delta, \mu, A}(\sigma(t_i')) = \sigma_{\mathcal{I}_{\Delta, \mu, A}(t_i')}\). This implies that \(f(t_1', \ldots, t_k') = \sigma_{\mathcal{I}_{\Delta, \mu, A}(t)}\).

Lemma 3. Let \(\tau = (P, R, S, \mu)\) be a CS problem where \(P\) and \(R\) are TRSs over the signatures \(G\) and \(F\) respectively and \(\Delta \subseteq F\). Let \(A\) be an infinite strongly minimal \((P, R, S, \mu)\)-chain. Let \(P_G \cup U_C(R)\) be conservative and left-linear, \(P_X\) left-linear and conservative with respect to \(S\), and \(\Delta = \{\text{root}(\ell) \mid \ell \rightarrow r \in U_C(R)\}\). If \(s\) and \(t\) are \((R, \mu)\)-terminating and \(s \Rightarrow_{\mathcal{R}, \mu} t\) then \(\mathcal{I}_{\Delta, \mu, A}(s) \Rightarrow_{\mathcal{I}_{\Delta, \mu, A}(t)} \mathcal{I}_{\Delta, \mu, A}(s) \Rightarrow_{\mathcal{C}_t, \mu} \mathcal{I}_{\Delta, \mu, A}(t)\).

Proof. We proceed by induction on the position \(p\) of the redex in \(s \Rightarrow_{(\ell \rightarrow r), \mu} t\).

First assume that \(\text{root}(s) \in \Delta\) and \(p = \varepsilon\) (and therefore \(\ell \rightarrow r \in U_C(R)\)). So we have \(s = \sigma(\ell) \Rightarrow_{\ell \rightarrow r) \mu} \sigma(r) = t\). Moreover, for all subterms \(r'\) at active positions of \(r\), \(\text{root}(r')\) is \(\Delta\) and for all subterms \(r''\) at frozen positions of \(r\), \(\text{root}(r'')\) is \(\Delta\) by definition of \(\Delta\). We know that \(\ell \rightarrow r\) is conservative and left-linear. We have:

\[
\mathcal{I}_{\Delta, \mu, A}(s) = \mathcal{I}_{\Delta, \mu, A}(\sigma(\ell)) \Rightarrow_{\mathcal{C}_t, \mu} \sigma_{\mathcal{I}_{\Delta, \mu, A}(\ell)} \Rightarrow_{\ell \rightarrow r} \sigma_{\mathcal{I}_{\Delta, \mu, A}(\ell)} = \mathcal{I}_{\Delta, \mu, A}(\sigma(r)) = \mathcal{I}_{\Delta, \mu, A}(t)
\]

by Lemma 2.

Note that, by conservativity, every variable at an active position in \(r\) is at an active position in \(\ell\). Now consider the case where \(\text{root}(s) \in \Delta\) and \(p \neq \varepsilon\). Hence, \(s = f(s_1, \ldots, s_k), t = f(s_1', \ldots, s_k')\), \(i \in \mu\)(f), and \(s_i' \Rightarrow_{(\ell \rightarrow r), \mu} t_i\).
The induction hypothesis implies \( I_{\Delta,\mu,A}(s') \rightarrow^{+}_{t \rightarrow r} \sigma \in C_{\tau,\mu} \ I_{\Delta,\mu,A}(t) \) and hence, \( I_{\Delta,\mu,A}(s) \rightarrow^{+}_{t \rightarrow r} \sigma \in C_{\tau,\mu} \ I_{\Delta,\mu,A}(t) \). Finally, we consider the case \( \text{root}(s) \notin \Delta \). In this case, \( I_{\Delta,\mu,A}(s) \in \text{order} \{ \{ I_{\Delta,\mu,A}(u) \mid \ s \rightarrow_{\mathcal{T},\mu} u \} \} \) because \( s \rightarrow_{\mathcal{T},\mu} t \). By applying \( C_{\tau} \) rules, we get \( I_{\Delta,\mu,A}(s) \rightarrow^{+}_{t \rightarrow r} \sigma \in C_{\tau,\mu} \ I_{\Delta,\mu,A}(t) \).

**Proof (Theorem 3).** Regarding soundness, we proceed by contradiction. By Theorem 2, assume that there is an infinite strongly minimal \( (\mathcal{P}, R, S, \mu) \)-chain \( A \), but there is no infinite strongly minimal \( (\mathcal{P} - P_{\neg}, R, S - S_{\neg}, \mu) \)-chain. Due to the finiteness of \( \mathcal{P} \) and \( S \), we can assume that there are subsets \( Q \subseteq \mathcal{P} \) and \( T \subseteq S \) such that \( A \) has a tail \( B \)

\[
\sigma(u_1) \left\{ \begin{array}{l}
\rightarrow_{Q_X} \circ \xi_{\mathcal{T},\mu} \circ \xi_T \\
\rightarrow_{Q_X} \circ \xi_{\mathcal{T},\mu} \circ \xi_T
\end{array} \right. \rightarrow_{N,\mu} \sigma(u_2) \left\{ \begin{array}{l}
\rightarrow_{Q_X} \circ \xi_{\mathcal{T},\mu} \circ \xi_T \\
\rightarrow_{Q_X} \circ \xi_{\mathcal{T},\mu} \circ \xi_T
\end{array} \right. \cdots
\]

for some substitution \( \sigma \), where all pairs in \( Q \) and all rules in \( T \) are infinitely often used (note that, if \( T \neq \emptyset \), then \( T \neq \emptyset \) and \( Q_X \neq \emptyset \)), and, for all \( i \geq 1 \), (1) if \( u_i \rightarrow v_i \in Q_G \), then \( v_i' = \sigma(v_i) \) and (2) if \( u_i \rightarrow v_i = u_i \rightarrow x_i \in Q_X \), then \( \sigma(x_i) \rightarrow_{\mathcal{T},\mu} \circ \xi_T \rightarrow_{T} t_i' \). Moreover, all \( t_i' \) are \((R,\mu)\)-terminating.

We know that \( \mathcal{P} \cup U(\mathcal{R}) \cup S_{\mathcal{P},\mu} \) is conservative and left-linear. We apply \( I_{\Delta,\mu,A} \) in Definition 7 to the initial term. To ease readability, we let \( I = I_{\Delta,\mu,A} \). By definition, we know that \( I(\sigma(u_1)) = I_{0,\Delta,\mu}(\sigma(u_1)) \). Let \( \Delta = \{ \text{root}(t) \mid t \rightarrow r \in U(\mathcal{R}) \} \). Note that the application of \( I_{0,\Delta,\mu} \) is always possible since active subterms are \((R,\mu)\)-terminating due to the minimality of the chain. Using Lemma 3, we obtain \( I(\sigma(u_1)) \rightarrow^{*}_{U(\mathcal{R}) \cup C_{\tau,\mu}} I(\sigma(u_{i+1})) \) for all \( i \geq 1 \). Moreover, by the definition of \( U(\mathcal{R}) \), for all non-variable subterms \( v_i' \) at active positions of \( v_i \), we have \( \text{root}(v_i') \in \Delta \). By Lemma 3, we have \( \sigma_i(v_i) = I(\sigma(v_i)) \) and \( I(\sigma(u_{i+1})) \rightarrow^{*}_{C_{\tau,\mu}} \sigma_i(u_{i+1}) \). Since \( u_i \) and \( \sigma_i(v_i) \) for all \( u_i \rightarrow v_i \in Q \subseteq \mathcal{P} \), by stability of \( \mathcal{R} \) and \( \mathcal{P} \), we have \( \sigma_i(u_i) \sim (\mathcal{R} \cup \mathcal{P}) \sigma_i(v_i) \) for all \( i \geq 1 \).

No pair \( u \rightarrow v \in Q \) satisfies that \( u \sim v \). Otherwise, we get a contradiction by considering the following two cases:

1. If \( u_i \rightarrow v_i \in Q_{\mathcal{R}} \), then \( \sigma_i(v_i) \rightarrow^{*}_{U(\mathcal{R}) \cup C_{\tau,\mu}} \sigma_i(u_{i+1}) \) and (by compatibility of \( \mathcal{R} \) with the rules in \( U(\mathcal{R}) \cup C_{\tau,\mu} \)) \( \sigma_i(u_i) \) and \( \sigma_i(v_i) = I(\sigma(v_i)) \) (we use the fact that \( \mathcal{P} \) is conservative and left-linear). Since \( \sigma_i(u_i) \sim (\mathcal{R} \cup \mathcal{P}) \sigma_i(v_i) \), by transitivity of \( \mathcal{R} \) and compatibility between \( \mathcal{R} \) and \( \mathcal{P} \), we conclude that \( \sigma_i(u_i) \sim (\mathcal{R} \cup \mathcal{P}) \sigma_i(u_{i+1}) \).

2. If \( u_i \rightarrow v_i = u_i \rightarrow x_i \in Q_X \) (which is not empty whenever \( T \) is not empty), then \( \sigma_i(v_i) = \sigma_i(x_i) \rightarrow^{*}_{T_{\mathcal{P},\mu \cup C_{\tau,\mu}}} I(s_i) \), as in Lemma 3. Let \( \ell_i \rightarrow r_i \in T_{\mathcal{R}} \).

For all subterms \( r_i' \) at active positions of \( r_i \), \( \text{root}(r_i') \in \Delta \). If \( \ell_i \) is not linear, we know that \( u_i \neq \sigma(\ell_i) \) and there is a rule \( u \rightarrow v \in \mathcal{P} \cup \mathcal{R} \) such that \( \sigma(u) \neq \sigma(\ell_i) \) and \( \sigma(v) \rightarrow_{\mathcal{P},\mu} \sigma(\ell_i) \) (if \( \sigma(v) \rightarrow_{\mathcal{P},\mu} \sigma(\ell_i) \), we know that \( \sigma(\ell_i) \) is not \((R,\mu)\)-terminating, that violates a premise of the paper) and we know that \( u \) is left-linear and \( v = \ell_i \) up to renaming of variables, so \( \sigma_i \) is unique. Since \( \sigma_i(v_i) \sim (\mathcal{R} \cup \mathcal{P}) \sigma_i(v_i) \) for all \( \ell \rightarrow r \in T \), we have \( \sigma_i(v_i) = \sigma_i(x_i) \sim (\mathcal{R} \cup \mathcal{P}) I(s_i) \sim (\mathcal{R} \cup \mathcal{P}) \sigma_i(r_i) = I(t_i) \). Hence, by transitivity of \( \mathcal{R} \) (and compatibility of \( \mathcal{R} \) and \( \mathcal{P} \)),
Theorem 3. The following polynomial interpretation removes pair

\[ \text{(17)} \]

The following TRS

Example 5.

and we obtain a finite CS problem.

\( \tau \)

The new CS problem

numbers using CSR

\( \mu \)

Since \( u \to v \) occurs infinitely often in \( B \), there is an infinite set \( \mathcal{I} \subseteq \mathbb{N} \) of pairs such that \( \sigma_I(u_i) \not\approx \sigma_I(u_{i+1}) \) for all \( i \in \mathcal{I} \). Thus, by using the compatibility conditions of the \( \mu \)-reduction pair, we obtain an infinite decreasing \( \not\approx \)-sequence which contradicts well-foundedness of \( \not\approx \). Therefore, \( B \) is an infinite minimal \( (\mathcal{P} - \mathcal{P}, \mathcal{R}, \mathcal{S} - \mathcal{S}, \mu) \)-chain, thus leading to a contradiction.

6 Application Examples

Example 4. Continuing with Example 1, and the CS problem presented at the beginning of Section 3, \( \tau = \{(15), (17), (18)\}, \mathcal{R}, \mathcal{S}, \mu \), where \( \mathcal{S} = \{x+y \to y, \text{fact}(p(x)) \to \text{FACT}(p(x))\} \). Up to now, no automatic tool was able to handle this CS problem. By Definition 5, the set of extended basic CS usable rules \( \mathcal{U}_\tau^e(\mathcal{R}) \) is:

\[
\begin{align*}
\text{zero}(0) & \to \text{true} & \text{zero}(s(x)) & \to \text{false} & p(s(x)) & \to x
\end{align*}
\]

when previously all the rules are usable. We can use the extended basic CS usable rules instead of \( \mathcal{R} \) because the CS problem satisfies the restrictions in Theorem 3 and the following polynomial interpretation [19] allows us to remove pair (18):

\[
\begin{align*}
\text{fact}(x) & = 2x & *x, y & = \frac{1}{2}xy & p(x) & = \frac{1}{2}x & \text{zero} & = \frac{1}{2}x^2 \\
0 & = 2 & s(x) & = x + 1 & \text{false} & = 2 & \text{true} & = 2 \\
\text{FACT}(x) & = 2x^2 & \text{IF}(x, y, z) & = \frac{1}{2}xy + \frac{1}{2}x + z
\end{align*}
\]

The new CS problem \( \tau' = \{(15), (17)\}, \mathcal{R}, \mathcal{S}, \mu \) can be handled again using Theorem 3. The following polynomial interpretation removes pair (17):

\[
\begin{align*}
\text{fact}(x) & = 1 & *x, y & = 2 & p(x) & = 2x + 1 & \text{zero} & = 2x + 1 \\
0 & = 0 & s(x) & = x & \text{false} & = 1 & \text{true} & = 1 \\
\text{FACT}(x) & = 2 & \text{IF}(x, y, z) & = 2y + 1
\end{align*}
\]

and we obtain a finite CS problem.

Example 5. The following TRS \( \mathcal{R} \) [10] models the subtraction of two natural numbers using CSR and if-then-else:

\[
\begin{align*}
0 \leq y & \to \text{true} & s(x) \leq 0 & \to \text{false} \\
s(x) \leq s(y) & \to x \leq y & p(0) & \to 0 \\
p(s(x)) & \to x & x - y & \to \text{if}(x \leq y, 0, s(p(x) - y)) \\
\text{if(true, } x, y) & \to x & \text{if(false, } x, y) & \to y
\end{align*}
\]

with \( \mu(-) = \mu(\leq) = \{1, 2\} \), \( \mu(s) = \mu(\text{if}) = \mu(p) = \{1\} \) and \( \mu(0) = \mu(\text{true}) = \mu(\text{false}) = 0 \). Without the replacement map, this example is non-terminating.
Consider the following CS problem \( \tau = (\mathcal{P}, \mathcal{R}, \mathcal{S}, \mu) \) where \( \mathcal{P} \) is:

\[
\begin{align*}
  x - y &\rightarrow \text{IF}(x \leq y, 0, s(p(x) - y)) \quad (22) \\
  \text{IF}(\text{false}, x, y) &\rightarrow y \\
  \text{IF}(\text{true}, x, y) &\rightarrow x \\
  \text{IF}(\text{false}, x, y) &\rightarrow y
\end{align*}
\]

with \( \mu \) extended by \( \mu(-^) = \{1, 2\} \), and \( \mu(\text{IF}) = \{1\} \), where \( \mathcal{S} = \{ s(x) \rightarrow x, p(x) - y \rightarrow p(x) - ^y \} \). By Definition 5, the set of extended basic CS usable rules \( \mathcal{U}_e^{\tau}(\mathcal{R}) \) is:

\[
\begin{align*}
  0 \leq y &\rightarrow \text{true} \\
  s(x) \leq 0 &\rightarrow \text{false} \\
  s(x) \leq s(y) &\rightarrow x \leq y \\
  p(0) &\rightarrow 0 \\
  p(s(x)) &\rightarrow x
\end{align*}
\]

when previously all the rules are usable. We can use the extended basic CS usable rules instead of \( \mathcal{R} \) because the CS problem satisfies the restrictions in Theorem 3 and the following polynomial interpretation allows us to remove pair (24):

\[
\begin{align*}
  [\text{-}](x, y) &= \frac{1}{2} x + \frac{1}{2} y \\
  [\text{0}] &= 0 \\
  [\text{true}] &= 0 \\
  [\text{false}] &= 0 \\
  [\text{-}^y](x, y) &= x + y + 1
\end{align*}
\]

The new CS problem \( \tau' = (\{22\}, \mathcal{R}, \mathcal{S}, \mu) \) can be handled again using Theorem 3. The following polynomial interpretation removes pair (22):

\[
\begin{align*}
  [-](x, y) &= 2x + 2y + 2 \\
  [\text{0}] &= 0 \\
  [\text{true}] &= 0 \\
  [-^y](x, y) &= 2y + 2 \\
  [\text{IF}(x, y, z)] &= x + y + 1
\end{align*}
\]

and we obtain a finite CS problem.

7 Conclusions

In this paper, we analyze why dealing with if-then-else expressions in CSR is hard in practice and we propose a new notion of chain, the notion of strongly minimal \((\mathcal{P}, \mathcal{R}, \mathcal{S}, \mu)\)-chain and a new set of CS usable rules, the extended basic usable rules, that allows us to simplify termination proofs involving if-then-else expressions. Furthermore, we show a couple of examples (included in the TPDB) where the technique is successfully applied: Example 1 in [22] and Example 1 in [10], whose automatic proofs are an open problems since 1997 and 2003.

As future work, we propose to study these definitions in connection to innermost termination and to implement the new definitions in our tool for proving termination properties, MU-TERM [5].

References

Partial Evaluation for Java Malware Detection

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Abstract. The fact that Java is platform independent gives hackers the opportunity to write exploits that can target users on any platform, which has a JVM implementation. To circumvent detection by anti virus (AV) software, obfuscation techniques are routinely applied to make an exploit more difficult to recognise. Popular obfuscation techniques for Java include string obfuscation and applying reflection to hide method calls; two techniques that can either be used together or independently. This paper shows how to apply partial evaluation to remove these obfuscations and thereby improve AV matching. The paper presents a partial evaluator for Jimple, which is a typed three-address code suitable for optimisation and program analysis, and also demonstrates how the residual Jimple code, when transformed back into Java, improves the detection rates of a number of commercial AV products.

1 Introduction

Java is both portable and architecture-neutral. It is portable because Java code is compiled to JVM byte code for which interpreters exist, not only for the popular desktop operating systems, but for phones and tablets, and as browser plug-ins. It is architecture-neutral because the JVM code runs the same regardless of environment. This presents a huge advantage over languages, such as C/C++, but also poses a major security threat. If an exploit leverers a vulnerability in a JVM implementation, it will affect all versions of a JVM that have not closed off that loophole, and well as those users who have not updated their JVM.

JVM vulnerabilities have been used increasingly by criminals in so-called client side attacks, often in conjunction with social engineering tactics. For example, a client-side attack might involve sending a pdf document \cite{pdf} that is designed to trigger a vulnerability when it is opened by the user in a pdf reader. Alternatively a user might be sent a link to a website which contains a Java applet which exploits a JVM vulnerability \cite{applet} to access the user’s machine. Client-side attacks provide a way of bypassing a firewall that block ports to users’ machines and, are proving to be increasingly popular: last year many JVM exploits were added to the Metasploit package, which is a well-known and widely-used penetration testing platform. This, itself, exacerbates the problem. As well as serving penetration testers and security engineers, a script kiddie and or a skilled black-hat can reuse a JVM vulnerability reported in Metasploit, applying obfuscation so that it is not recognised by even up-to-date AV detection software.

Experimental evidence suggests that commercial AV software use Metasploit as source of popular attack vectors, since exploits from Metasploit are
typically detected if they come in an unadulterated form. One can only speculate what techniques an AV vendor actually uses, but detection methods range from entirely static techniques, such as signature matching, to entirely dynamic techniques, in which the execution of the program or script is monitored for suspicious activity. In signature matching, a signature (a hash) is derived, often by decompiling a sample, which is compared against a database of signatures constructed from known malware. Signatures are manually designed to not trigger a false positive which would otherwise quarantine an innocuous file. Dynamic techniques might analyse for common viral activities such as file overwrites and attempts to hide the existence of suspicious files, though it must be said, there are very few academic works that address the classification of Java applets [17].

The essential difference between running a program in an interpreter and partially evaluating it within a partial evaluator is that the latter operates over a partial store in which not all variables have known values; the store determines which parts of the program are executed and which parts are retained in the so-called residual. Partial evaluation warrants consideration in malware detection because it offers a continuum between the entirely static and the entirely dynamic approaches. In particular, one can selectively execute parts of the program, namely those parts that mask suspicious activity, and then use the residual in AV scanning. This avoids the overhead of full execution while finessing the problem of incomplete data that arises when a program is evaluated without complete knowledge of its environment. On the theoretical side, partial evaluation provides nomenclature (e.g., polyvariance) and techniques (e.g., generalisation) for controlling evaluation and specialisation. On the practical side, partial evaluation seems to be partially appropriate for AV matching because Java exploits are often obfuscated by string obfuscation and by using advanced language features such as reflection. Although reflection is designed for such applications as development environments, debuggers and test harnesses, it can also be applied to hide a method call that is characteristic of the exploit. This paper will investigate how partial evaluation can be used to deobfuscate malicious Java software; it argues that AV scanning can be improved by matching on the residual JVM code, rather than original JVM code itself.

1.1 Contributions

This paper describes how partial evaluation can deobfuscate malicious Java exploits; it revisits partial evaluation from the perspective of Java malware detection which, to our knowledge, is novel. The main contributions are as follows:

- The paper describes the semantics of a partial evaluator for Jimple [20]. Jimple is a typed three-address intermediate representation that is designed to support program analysis and, conveniently, can be decompiled back into Java for AV matching.
- The paper shows how partial evaluation can be used to remove reflection from Jimple code, as well as superfluous string operations, that can be used to obfuscate malicious Java code.
The paper describes how partial evaluation can be used in tandem with an abstract domain so as to avoid raising an error when a branch condition cannot be statically resolved [18, section 3.3]. In such a situation, one branch might lead to an environment whose bindings are inconsistent with that generated along another. Rather than abort when inconsistent environments are detected at a point of confluence, we merge the environments into an abstract environment that preserves information from both, so that partial evaluation can continue.

2 Primer on Java Obfuscation

This section will describe techniques that are commonly used to obfuscate Java code to avoid AV detection. The obfuscations detailed below are typically used in combination; it is not as if one obfuscation is more important than another.

2.1 Reflection Obfuscation

An AV filter might check for the invocation of a known vulnerable library function, and to thwart this, malicious applets frequently use reflection to invoke vulnerable methods. This is illustrated by the code in listing 1.1 which uses the `Class.forName` static method to generate an object `c` of type `Class`. The `c` object allows the programmer to access information pertaining to the Java class `java.security.Permissions`, and in particular create an object `o` of type `java.security.Permissions`. Furthermore, `c` can be used to create an object `m` that encapsulates the details of a method call on object `o`. The invocation is finally realised by applying the `invoke` on `m` using `o` as a parameter. This sequence of reflective operation serves to disguise what would otherwise be a direct call to the method `add` on an object of type `Permissions`.

2.2 String Obfuscation

Malicious applets will often assemble a string at run-time from a series of component strings. Alternatively a string can be encoded and then decoded at run-time. Either tactic will conceal a string, making it more difficult to recognise class and
method names, and thereby improving the chances of outwitting a signature-based AV system. Listing 1.2 gives an example of a string reconstruction method that we found in the wild, in which the string `java.lang.SecurityManager` is packed with numeric characters which are subsequently removed at runtime. Listing 1.3 illustrates an encoder which replaces a letter with the letter 13 letters after it in the alphabet. The encoded strings are then decoded at run-time before they are used to create a handle of type `Class` that can, in turn, be used to instantiate `java.lang.SecurityManager` objects.

2.3 Other Obfuscations

There is also no reason why other obfuscations [6] cannot be used in combination with reflection and string obfuscation. Of these, one of the most prevalent is name obfuscation in which the names of the user-defined class and method names are substituted with fresh names. For example, the name `getStr` in Listing 1.2 might be replaced with a fresh identifier, so as to mask the invocation of a known decipher method.

3 Partial Evaluation

In this section we outline a partial evaluator for removing string obfuscation and reflection from Jimple code, which is a three address intermediate representation (IR) for the Java programming language and byte code. Jimple is supported by the Soot static analysis framework and, quite apart from its simplicity, Soot provides support for translating between Jimple and Java.

There are two approaches to partial evaluation: online and offline. In the online approach specialisation decisions are made on-the-fly, based on the values of
public static String rot13(String s) {
    StringBuffer sb = new StringBuffer();
    for (int i = 0; i < s.length(); i++) {
        char c = s.charAt(i);
        if (c >= 'a' && c <= 'm') c += 13;
        else if (c >= 'A' && c <= 'M') c += 13;
        else if (c >= 'n' && c <= 'z') c -= 13;
        else if (c >= 'N' && c <= 'Z') c -= 13;
        sb.append(c);
    }
    return sb.toString();
}

String str = "wnin.ynat.Frphevg1Znantre";
Class<?> c = Class.forName(rot13(str));

Listing 1.3. String obfuscation using the rot13 substitution cipher
\[
\mathcal{P}[\text{var } t; S](\rho, \sigma, o) = \\
v = \text{default}(t) \\
l = \text{allocate}(v) \\
\rho' = \{x \mapsto l\} \\
\sigma' = \{l \mapsto v\} \\
\text{emit}[\text{var } t] \\
\mathcal{P}[S](\rho \circ \rho', \sigma \circ \sigma', o)
\]

\[
\mathcal{P}[x := y \oplus z; S](\rho, \sigma, o) = \\
\text{if } \sigma(\rho(y)) \lor \sigma(\rho(z)) = \top \text{ then} \\
\text{emit}[x := y \oplus z] \\
\sigma' = \sigma \circ \{\rho(x) \mapsto \top\} \\
\mathcal{P}[S](\rho, \sigma', o) \\
\text{else} \\
v = \sigma(\rho(y)) \lor \sigma(\rho(z)) \\
\text{emit}[x := v] \\
\sigma' = \sigma \circ \{\rho(x) \mapsto v\} \\
\mathcal{P}[S](\rho, \sigma', o) \\
\text{endif}
\]

\[
\mathcal{P}[x := @this; S](\rho, \sigma, o) = \\
\text{emit}[x := @this] \\
\sigma' = \{\rho(x) \mapsto o\} \\
\mathcal{P}[S](\rho, \sigma \circ \sigma', o)
\]

\[
\mathcal{P}[x := @\text{parameter}; S](\rho, \sigma, o) = \\
\text{emit}[x := @\text{parameter}] \\
\sigma' = \{\rho(x) \mapsto \sigma(\rho(\text{parameter}))\} \\
\mathcal{P}[S](\rho, \sigma \circ \sigma', o)
\]

\[
\mathcal{P}[x := \text{new } C; S](\rho, \sigma, o) = \\
\langle l, f \rangle = \text{getFields}(C) \\
l = \text{allocates}(f) \\
\bar{v} = \text{defaults}(l) \\
\rho' = \{f_0 \mapsto l_0, \ldots, f_n \mapsto l_n\} \\
\sigma' = \{l_0 \mapsto v_0, \ldots, l_n \mapsto v_n, \rho(x) \mapsto C : \rho'\} \\
\text{emit}[x := \text{new } C] \\
\mathcal{P}[S](\rho, \sigma \circ \sigma', o)
\]

\textbf{Fig. 1.} Outline of partial evaluator: declarations and assignments
\[ P[\text{return } x; \{p, \sigma, o\}] = \]
\[ \sigma' = \{p(\text{return}) \mapsto \sigma(p(x))\} \]
\[ \text{emit}[\text{return } x] \]
\[ \langle p, \sigma \circ \sigma', o \rangle \]

\[ P[x := \text{virtualinvoke}(obj, m(\bar{t}), \bar{y}); S]\{p, \sigma, o\} = \]
\[ \text{if } \sigma(p(obj)) = \top \text{ then} \]
\[ \text{emit}[x := \text{virtualinvoke}(obj, m(\bar{t}), \bar{y})] \]
\[ \sigma' = \{p(x) \mapsto \top\} \]
\[ P[S]\{p, \sigma \circ \sigma', o\} \]
\[ \text{else if } \sigma(p(obj)) = C : \rho' \]
\[ \text{if } C = \text{Method} \land m = \text{invoke} \]
\[ \text{if } \sigma(\rho(\text{method})) = \text{null} \text{ then error} \]
\[ \bar{y}' = \langle y_1, \ldots, y_n \rangle \]
\[ \text{emit}[x := \text{virtualinvoke}(y_0, \sigma(\rho(\text{method})), \bar{y}')] \]
\[ \sigma' = \{p(x) \mapsto \top\} \]
\[ P[S]\{p, \sigma \circ \sigma', o\} \]
\[ \text{else} \]
\[ \bar{v} = \sigma(p(\bar{y})) \]
\[ B = \text{findMethod}(C.m(\bar{t})) \]
\[ \bar{l} = \text{allocates}(\bar{v}) \]
\[ k = \text{allocate}(\text{result}) \]
\[ \rho'' = \{\text{parameter}_0 \mapsto l_0, \ldots, \text{parameter}_n \mapsto l_n, \text{result} \mapsto k\} \]
\[ \sigma' = \sigma \circ \{l_0 \mapsto v_0, \ldots, l_n \mapsto v_n\} \]
\[ \langle \sigma'' \rangle = P[B]\{\rho'' \circ \rho', \sigma', C : \rho'\} \]
\[ \sigma''' = \{p(x) \mapsto \sigma''(p'(\text{result}))\} \]
\[ P[S]\{p, \sigma \circ \sigma''', o\} \]

\[ P[\text{if } x \text{ goto } l; S]\{p, \sigma, o\} = \]
\[ \text{if } \sigma(p(x)) = 0 \text{ then} \]
\[ P[S]\{p, \sigma, o\} \]
\[ \text{else if } \sigma(p(y)) = 1 \text{ then} \]
\[ S = \text{lookup}(l) \]
\[ P[S]\{p, \sigma, o\} \]
\[ \text{else if } \sigma(p(x)) = \top \text{ then} \]
\[ n = \text{lookup}(S, P) \]
\[ c = \text{confluence}(l, n) \]
\[ \text{emit}[\text{if } n \text{ goto } l] \]
\[ T_i = \text{branch}(l, c) \]
\[ T_f = \text{branch}(n, c) \]
\[ \langle p, \sigma_i, o \rangle = P[T_i]\{p, \sigma, o\} \]
\[ \langle p, \sigma_f, o \rangle = P[T_f]\{p, \sigma, o\} \]
\[ \sigma' = \{l \mapsto v \mid l \in \text{codomain}(\rho) \land v = \text{if } \sigma_i(l) = \sigma_f(l) \text{ then } \sigma_i(l) \text{ else } \top\} \]
\[ P[P \text{ drop } c]\{p, \sigma', o\} \]

Fig. 2. Outline of partial evaluator: control-flow
3.1 type declarations

A statement \( \text{var } t \ x \) declares that \( x \) is of type \( t \), where \( t \) is either primitive or a user-defined class. Such a declaration is handled by allocating a memory location \( l \) using the auxiliary \( \text{allocate} \) and then updating the environment \( \rho' \) to reflect this change. The store is also mutated to map location \( l \) to the default value for the type \( t \), which is given by the auxiliary function \( \text{default} \). The default values for the primitives types are 0 for \( \text{int} \) and 0 for \( \text{boolean} \). The default types for object types is \( \text{null} \).

3.2 new

A statement \( x = \text{new } C \) instantiates the class \( C \) to create an object that is represented by a pair \( C : \rho \) where the environment \( \rho \) maps the fields of \( C \) to memory locations that store their values. Different objects \( C : \rho_1 \) and \( C : \rho_2 \) from the same class \( C \) map the same field variables to different locations. The auxiliary method \( \text{getFields} \) retrieves the types and the names of the fields of the class \( C \). The function \( \text{defaults} \) takes a vector of types \( t \) and returns a vector of default values that is used to populate the fields, following the conventions of \( \text{default} \).

3.3 arithmetical operations

An assignment statement \( x := y \oplus z \) can only be statically evaluated if both the variables \( y \) and \( z \) are bound to known values. In this circumstance the assignment \( x := y \oplus z \) is specialised to \( x := v \) where \( v \) is the value of the expression \( y \oplus z \). The store \( \sigma' \) is updated to reflect the new value of \( x \), as the value of \( x \) is statically known. Note that the residual includes \( x := v \), even though information on \( x \) is duplicated in the store \( \sigma' \), so as to permit subsequent statements, with reference \( x \), to be placed in the residual without substituting \( x \) with its value \( v \). If there are no statements that reference \( x \) then the assignment \( x := v \) will be removed by dead variable elimination, which is applied as a post-processing step.

3.4 this and parameters

In Jimple there is a distinguished variable \( \text{this} \) which stores the current object reference which, in the partial evaluator, is modelled with the current object \( o \), that is passed with the environment and the store. An assignment statement \( x := @\text{this} \) thus merely updates the location \( \rho(x) \) with \( o \).

Also in Jimple, a special variable \( \text{parameter}_i \) is used to store the location of the \( i \)th formal argument of a method call, where the first argument has an index of 0. This is modelled more directly in the partial evaluator, so that an assignment statement \( x := @\text{parameter}_i \) updates the location \( \rho(x) \) with the value of this parameter.
3.5 return and virtualinvoke

The statement \textit{return} \(x\) writes the value of \(x\) to a special variable \textit{return}, which is subsequently read by \textit{virtualinvoke}.

The handling of \textit{virtualinvoke}(\textit{obj}, \textit{m}(\textit{t}), \textit{y}) is worthy of special note, both in the way reflective and non-reflective calls are handled. A reflective call arises when the method \(m\) coincides with \textit{invoke} and the object \textit{obj} is of type \textit{Method}. The reflective method call is a proxy for the so-called reflected method call. The reflected method call is not applied to the object \textit{obj} but an object that is prescribed by the first parameter of \(y\). Moreover, the method that is applied to this object is given by \textit{obj} in a distinguished field that, for our purposes, is called \textit{method}. This field either contains \texttt{null}, indicating that it has been initialised but not been reset, or a string that represents the name of a method that is to be invoked. If the method field is \texttt{null} then an error is issued, otherwise the string stored in the field \textit{method} is used to generate the residual. Note that the reflected method call is not invoked; merely placed in the residual. Note too that the first argument of \textit{virtualinvoke} in the residual is \(y_0\) whereas the last is the vector \(y'\) which coincides with \(y\) with the exception that the first element has been removed. The first argument is the variable name (rather than the object itself) to which the residuated method will ultimately be applied; the third argument is the list of actual arguments that will be passed to the method on its invocation.

In the case of a non-reflective call, the values of the parameters are looked up, and then an auxiliary function \texttt{findMethod} is applied to find a block \(B\) which is the entry point into the method of the class \(C\) whose signature matches \(m(\textit{t})\). The function \texttt{allocates} is then called to generate fresh memory locations, one for each actual argument \(y_0, \ldots, y_n\). The environment is then extended to map the distinguished variables \texttt{parameter}_0, \ldots, \texttt{parameter}_n to fresh memory locations, so as to store the actual arguments. The partial evaluator is then recursively involved on the block \(B\) using \(C : \rho'\) as the object. The net effect of the method call is to update the store \(\sigma'\) which is used when evaluating the remaining statements \(S\).

Note that this formulation assumes that method calls are side-effect free. Although this is true for string obfuscation methods that we have seen, for full generality the partial evaluator should be augmented with an alias analysis, in the spirit of side-effect analysis \[9\], that identifies those variables that are possibly modified by a method call, hence must be reset to \(\top\).

3.6 goto

Specialising a conditional jump in Jimple is not conceptually difficult, but is complicated by the way if \(x\) \texttt{goto} \(l; S\) will drop through to execute the first instruction of the sequence \(S\) if the boolean variable \(x\) is false. This makes the control-flow more difficult to recover when the value of \(x\) is unknown. When \(x\) is known to be true, however, the partial evaluator is merely redirected at a block \(B\) that is obtained by looking up the sequence whose first instruction is labelled
by \( l \). Conversely, if \( x \) is known to be false, then partial evaluation immediately proceeds with \( S \).

In the case when \( x \) has an undetermined value, the partial evaluator explores both branches until the point of confluence when both branches merge. Then the partial evaluator continues at the merge point, relaxing the store to \( \sigma' \) so that it is consistent with the stores that are derived on both branches. Note that partial evaluation does not halt if the stores are inconsistent; instead it will unify the two stores by replacing any inconsistent assignment for any location with an assignment to \( \top \). Note that it is only necessary to unify those locations that are reachable from the variables that are currently in scope.

To realise this approach, an auxiliary function \texttt{lookup} is used to find the position \( n \) of the first statement of \( S \) in the list \( P \) which constitutes the statements for the currently executing method. This is the position of the first statement immediately after the conditional branch. Then a function \texttt{confluence} examines the control-flow graph of the method so as to locate the confluence point, of the true and false branches, identified by the index \( c \) of \( S \). Both branches are evaluated separately with two copies of the environment, until the confluence point where the two environments are merged. Partial evaluation then resumes at the confluence point, which corresponds to the instruction sequence \( P \texttt{drop} c \), namely, execution is continued at the \( c^{th} \) instruction of the sequence \( P \).

3.7 Example

Listing 1.4 gives the before and after for a method call that is obfuscated by reflection and string obfuscation, using the ROT13 simple letter substitution cipher given in Listing 1.3. The residual Jimple code is presented as Java for readability. Completely unfolding the \texttt{rot13} method call decrypts the string \texttt{qbFbzrguvat} as the string \texttt{doSomething}. This string statically defines the value of the object \texttt{method}, allowing \texttt{method.invoke(m, null)} to be specialised to \texttt{m.doSomething()}, thereby removing the reflective call. Note that the variables \texttt{encrypted} and \texttt{method} cannot be removed without dead variable elimination.

4 Experiments

To assess how partial evaluation can aid in AV matching, a number of known applet malware samples from the Metasploit exploit package [2] were obfuscated using the techniques outlined in section 2. Details of the samples are given in Fig. 3; the samples were chosen entirely at random. So as to assess the effect of partial evaluation against a representative AV tool, we compared the detection rates, with and without partial evaluation, on eight commercial AV products. Together these products cover the majority of the global market, as reported in 2013 [15] and is illustrated in the chart given in Figure 3. Conveniently, VirusTotal [19] provides a prepackaged interface for submitting malware samples to all of these products, with the exception of Avira, which is why this tool does not appear in our experiments.
Listing 1.4. Before and after partial evaluation

Unfortunately, developing a complete partial evaluator for Jimple is a major undertaking, since it is necessary to support the entire Java API and runtime environment, which itself is huge. To side-step this engineering effort, we implemented a partial evaluator in Scala, following the description in section 3, only providing functionality for String, StringBuffer and StringBuilder classes. This was achievable since Java String objects are accessible to Scala. (Scala’s parser combinator library also make it is straightforward to engineer a parser for Jimple.) Although other objects could be handled in the same way, we simply took each of these obfuscated CVEs and extracted the Jimple code and methods that manipulated strings. This code was then partially evaluated so as to deobfuscate the string handling. The CVEs were then hand-edited to reflect the residual, and then ran through VirusTotal to check that the effects of obfuscation had been truly annulled. Future implementation work will be to automate the entire process, namely translate the Jimple residual into Java using Soot and then invoke VirusTotal automatically through its public web API.

Table [1] details the detection rates for the AVs given in Figure 3, without obfuscation, with just string obfuscate, with just reflection obfuscation, and with both obfuscations applied. This gives four experiments in all. It is important to appreciate that the obfuscations used in the fourth experiment include all those obfuscations introduced in the second and third experiments and no more.

The results show that in most cases the AVs detect most of the exploits in their unadulterated form. Exploits CVE-2012-5088 and CVE-2013-2460 go
the most undetected, which is possibly because both exploits make extensive use of reflection. It is interesting to see that the product with the highest market share (Microsoft) was unable to detect any of the exploits after string obfuscation, which suggests the removing this obfuscation alone is truly worthwhile. Moreover, after introducing reflection the AV detection count for each exploit drops significantly. Furthermore, applying reflection with string obfuscation is strictly stronger than applying string obfuscation and reflection alone. CVE-2012-4681 represents an anomaly under McAfee since reflection obfuscation impedes detection whereas, bizarrely, using reflection with string obfuscation does not. Interestingly, McAfee classifies this CVE with the message Heuristic.BehavesLike.Java.Suspicious-Dldr.C which suggests that it is using a heuristic behavioural approach which might explain its unpredictability.

Most importantly, applying partial evaluation to the CVEs used in the fourth experiment restores the detection rates to those of the first experiment. Thus detection is improved, without having to redistribute the signature database.

5 Related Work

Although there has been much work in Java security, partial evaluation and reflection, there are few works that concern all three topics. This section provides pointers to the reader for the main works in each of these three separate areas.

One of the very few techniques that has addressed the problem of detecting malicious Java Applets is Jarhead [17]. This recent work uses machine learning to detect malicious Applets based on 42 features which include such things as the number of instructions, the number of functions per class and cyclomatic complexity [19]. Jarhead also uses special features that relate to string obfuscation, such as the number and average length of the strings, and the fraction of strings that contain non-ASCII printable characters. Other features that it applies determine the degree of active code obfuscation, such as the number of times that reflection is used within the code to instantiate objects and invoke methods. Out of a range of classifiers, decision trees are shown to be the most
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reliable. Our work likewise aspires to be static, though partial evaluation takes this notion to the limit, so as to improve detection rates. Moreover, machine learning introduces the possibility of false negatives and, possibly worse, false positives. Our approach is to scaffold off existing AV products that have been carefully crafted to not trigger false positives, and improve their matching rates by applying program specialisation as a preprocessing step.

The objective of partial evaluation is to remove interpretive overheads from programs. Reflection can be considered to be one such overhead and therefore it is perhaps not surprising that it has attracted interest in the static analysis community; indeed the performance benefit from removing reflection can be significant. Civet represents state-of-the-art in removing Java reflection; it does not apply binding-time analysis (BTA) but relies on programmer intervention, using annotation to delineate static from dynamic data, the correctness of which is checked at specialisation time. Advanced BTAs have been defined for specialising Java reflection, though to our knowledge, none have been implemented. We know of no partial evaluator for Jimple, though Soot represents the ideal environment for developing one. Quite apart from its role in deob-
fuscation, partial evaluation can also be applied in obfuscation \[10\]: a modified interpreter, that encapsulates an obfuscation technique, is partially evaluated with respect to the source program to automatically obfuscate the source. Program transformation has been proposed for deobfuscating binary programs \[5\], by unpacking and removing superfluous jumps and junk, again with the aim of improving AV scanning. This suggest that partial evaluation also has a role in binary analysis, where the aim is to make malware detection more semantic \[7\].

Reflection presents a challenge for program analysis: quite apart from writes to object fields, reflection can hide calls, and hence mask parts of the call-graph so that an analysis is unsound. Points-to analysis has been suggested \[12\] as a way of determining the targets of reflective calls which, in effect, traces the flow of strings through the program. This is sufficient for resolving many, but not all calls, hence points-to information is augmented with user-specified annotation so as to statically determine the complete call graph. The use of points-to information represents an advance over using dynamic instrumentation to harvest reflective calls \[11\] since instrumentation cannot guarantee complete coverage. Partial evaluation likewise traces the flow of strings through the program, though without refining points-to analysis, it is not clear that it has the precision to recover the targets of reflective calls that have been willfully obfuscated with such techniques as a substitution cipher (rot13).

## 6 Future Work

Termination analysis is a subfield of static analysis within itself and thus far we have not explored how termination can improve unfolding. We simply unfold loops where the loop bound is known at specialisation time. We also deliberately do not unfold recursive methods, though this is a somewhat draconian limitation. Future work will aim to quantify how termination analysis can be applied in an online setting to improve the quality of malware detection.

Although we have not observed this in the wild, there is no reason why reflection cannot be applied to a method that obfuscates a string, such as a decryptor. This would thwart our approach to deobfuscation since the reflected call would be deobfuscated in the residual, but would not actually be evaluated on a given string. Thus we will explore how partial evaluation can be repeatedly applied to handle these multi-layered forms of obfuscation.

We will also examine how partial evaluation can remove less common forms of Java obfuscation such as control flow obfuscation and serialization and deserialization obfuscation, the latter appearing to be as amenable to partial evaluation as string obfuscation. In the long term we will combine partial evaluation with code similarity matching, drawing on techniques from information retrieval.

## 7 Conclusion

We have presented a partial evaluator for removing string and reflection obfuscation from Java programs, with the aim of improving the detection of malicious
Java code. Our work puts partial evaluation in a new light: previous studies have majored on optimisation whereas we argue that partial evaluation has a role in anti-virus matching. To this end, a partial evaluator has been designed for Jimple, which was strength tested on five malware samples from the Metasploit exploit framework, obfuscated using string and reflection obfuscation.

References

Access control and obligations
in the category-based metamodel:
a rewrite-based semantics*

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Abstract. We define an extension of the category-based access control (CBAC) metamodel to accommodate a general notion of obligation. Since most of the well-known access control models are instances of the CBAC metamodel, we obtain a framework for the study of the interaction between authorisation and obligation, such that properties may be proven of the metamodel that apply to all instances of it. In particular, the extended CBAC metamodel allows security administrators to check whether a policy combining authorisations and obligations is consistent.

Key Words: Security Policies, Access Control, Obligations, Rewriting

1 Introduction

Access control policies specify which actions users are authorised to perform on protected resources. An authorisation may entail an obligation to perform another action on the same or another resource. Standard languages for the specification of access control policies include also a number of primitives to specify obligations associated with authorisations. For example, within XACML, an obligation is a directive from the Policy Decision Point (PDP) to the Policy Enforcement Point (PEP) specifying an action that must be carried out before or after an access is approved.

The notion of obligation helps bridge a gap between requirements and policy enforcement. For example, consider a hospital scenario in which any doctor may be authorised to read the medical record of a patient in an emergency situation, but in that case there is a requirement to inform the patient afterwards. Although access control models deal mainly with authorisations, incorporating the notion of an obligation facilitates the tasks of ensuring that obligations are enforced and ensuring that obligations are compatible with authorisations.

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A metamodel for access control, which can be specialised for domain-specific applications, has been proposed in [3]. It identifies a core set of principles of access control, abstracting away many of the complexities that are found in specific access control models, in order to simplify the tasks of policy writing and policy analysis. The metamodel focuses on the notion of a category, which is a class of entities that share some property. Classic types of groupings used in access control, like a role, a security clearance, a discrete measure of trust, etc., are particular instances of the more general notion of category. In category-based access control (CBAC), permissions are assigned to categories of users, rather than to individual users. Categories can be defined on the basis of user attributes, geographical constraints, resource attributes, etc. In this way, permissions can change in an autonomous way (e.g., when a user attribute changes), unlike, e.g., role-based models [1], which require the intervention of a security administrator.

In this paper, we define an extension of the CBAC metamodel to accommodate a general notion of an obligation, obtaining a formal framework for modelling and enforcing access control policies that involve authorisations and obligations. We show examples of application in the context of emergency management. We do not make any specific assumptions on the components of the system. Instead, we aim at defining an abstract model of access control and obligations that can be instantiated in various ways to satisfy specific requirements. To specify dynamic policies involving authorisations and obligations in the metamodel, we adjust the notion of an event used in previous work, and describe a set of core axioms for defining obligations.

Summarising, we provide an axiomatic definition of a generic framework for the specification of access control and obligation models, obtained by extending the CBAC metamodel with notions of obligations and duties; a rewrite-based operational semantics for the extended metamodel, dealing with authorisation and obligation assessment, including mechanisms for the resolution of conflicts between authorisations and obligations; and a rewrite-based technique to prove properties of access control policies involving obligations.

Overview: In Section 2, we recall the CBAC metamodel. Section 3 discusses obligations, Section 4 presents the extended CBAC metamodel, and Section 5 the operational semantics for obligations. In Section 6, we discuss related work, and in Section 7, conclusions are drawn and further work is suggested.

2 Preliminaries

We assume familiarity with basic notions on first-order logic and term-rewriting systems [2]. We briefly describe below the key concepts underlying the category-based metamodel of access control; see [3] for a detailed description.
Informally, a category is any of several distinct classes or groups to which entities may be assigned. Entities are denoted by constants in a many sorted domain of discourse, including: a countable set $C$ of categories, denoted $c_0, c_1, \ldots$, a countable set $P$ of principals, denoted $p_0, p_1, \ldots$ (we assume that principals that request access to resources are pre-authenticated), a countable set $A$ of named actions, denoted $a_0, a_1, \ldots$, a countable set $R$ of resource identifiers, denoted $r_0, r_1, \ldots$, a finite set $\mathcal{A}$ of possible answers to access requests (e.g., \{grant, deny, undetermined\}) and a countable set $\mathcal{S}$ of situational identifiers to denote environmental information. More generally, entities can be represented by a data structure (e.g., a principal could be represented by a term $\text{principal}(p_i, \text{attributeList})$), but constants will be sufficient for most examples in this paper. A permission is a pair $(a, r)$ of an action and a resource, and an authorisation is a triple $(p, a, r)$ that associates a permission with a principal.

The metamodel includes the following relations to formalise these notions:

- Principal-category assignment: $\text{PCA} \subseteq P \times C$, such that $(p, c) \in \text{PCA}$ iff a principal $p \in P$ is assigned to the category $c \in C$.

- Permission-category assignment: $\text{ARCA} \subseteq A \times R \times C$, such that $(a, r, c) \in \text{ARCA}$ iff action $a \in A$ on resource $r \in R$ can be performed by the principals assigned to the category $c \in C$.

- Authorisations: $\text{PAR} \subseteq P \times A \times R$, such that $(p, a, r) \in \text{PAR}$ iff a principal $p \in P$ can perform the action $a \in A$ on the resource $r \in R$.

**Definition 1 (Axioms).** The relation $\text{PAR}$ satisfies the following core axiom, where we assume that there exists a relationship $\subseteq$ between categories; this can simply be equality, set inclusion (the set of principals assigned to $c \in C$ is a subset of the set of principals assigned to $c' \in C$), or an application specific relation.

$$(a1) \forall p \in P, \forall a \in A, \forall r \in R,$$

$$(\exists c, c' \in C, (p, c) \in \text{PCA} \land c \subseteq c' \land (a, r, c') \in \text{ARCA}) \iff (p, a, r) \in \text{PAR}$$

Operationally, axiom $(a1)$ can be realised through a set of functions, as shown in [5]. We recall the definition of the function $\text{par}(P, A, R)$ below; it relies on functions $\text{pca}$, which returns the list of categories assigned to a principal, and $\text{arca}$, which returns a list of permissions assigned to a category.

$$(a2) \text{par}(P, A, R) \rightarrow \text{if } (A, R) \in \text{arca}^*(\text{below}(\text{pca}(P))) \text{ then } \text{grant} \text{ else } \text{deny}$$

The function $\text{below}$ computes the set of categories that are subordinate to any of the categories in the list $\text{pca}(P)$. The function $\in$ is a membership operator on lists, $\text{grant}$ and $\text{deny}$ are answers, and $\text{arca}^*$ generalises the function $\text{arca}$ to take into account lists of categories.
3 Obligations and Events

Obligations differ from permissions in the sense that, although permissions can be issued but not used, an obligation usually is associated with some mandatory action, which must be performed at a time defined by some temporal constraints or by the occurrence of an event. Therefore fulfilling the obligations may depend on the assignment of certain permissions. This can lead to undesirable interactions between permissions and obligations, where the obligations in a policy cannot be fulfilled given the set of assigned permissions in the system at a certain state. If obligations go unfulfilled (that is, become violated), this raises the question of accountability, that is, to whom shall be imputed responsibility for unfulfilled obligations. To take into account the above mentioned features, we will define obligations consisting of an action (an atomic or a composed action) on a resource to be performed during an interval specified by an initial event and a terminal event, and we will extend the CBAC metamodel to specify the assignment of obligations to principals and to study the interaction between obligations and permissions.

An event represents an action that happened in the system modelled. To define obligations associated with events, here we distinguish between event types, denoted by \( ge_i \) (e.g., registration for a course, a fire alarm, etc.) and specific events, denoted by \( e_i \) (e.g., the student Max M. registered for the Logic course in September 2012, the fire alarm went off in the Strand Building at 5pm on the 30 June 2012). A typing function will be used to classify events.

**Definition 2 (Event History and Interval).** An event history, denoted by \( h \in H \), is a sequence of events that may happen in a run of the system being modelled. A subsequence \( i \) of \( h \) is called an event interval; the first event in \( i \) opens the interval and the last one closes it.

We assume that an event history contains all the relevant events in order to determine, at each point, the set of authorisations and obligations of each principal (i.e., to determine the system state). Events and event types will be specified using a special-purpose language described in the next section. We assume that the typing relation associating event types with events is decidable.

Models of access control specify authorisations by defining the way permissions are assigned to principals. An obligation model should specify the way obligations are assigned to principals, but in addition, it should be possible to determine, at each point in the history, which obligations assigned to principals have to be enforced. For the latter, we introduce the notion of a duty.

**Definition 3 (Obligation).** A generic obligation is a tuple \( (a, r, ge_1, ge_2) \), where \( a \) is an action, \( r \) a resource, and \( ge_1, ge_2 \) two event types. If there is no starting
event (resp., no ending event) we write \((a, r, \bot, ge)\) (resp., \((a, r, ge, \bot)\)), meaning that the action on the resource must be performed at any point before an event of type \(ge\) (resp. at any point after an event of type \(ge\)).

**Example 1.** Assume that in an organisation, the members of the security team must call the fire-department if a fire alarm is activated, and this must be done before they de-activate the alarm. This obligation could be represented by the tuple \((\text{call, firedept, alarmON, alarmOFF})\);

**Definition 4 (Duty).** A duty is a tuple \((p, a, r, e_1, e_2, h)\), where \(p\) is a principal, \(a\) an action, \(r\) a resource, \(e_1, e_2\) are two events and \(h\) is an event history that includes an interval opened by \(e_1\) and closed by \(e_2\).

Unlike access control models, which do not need to check whether the authorised actions are performed or not by the principals, obligation models need to include mechanisms to check whether duties were discharged or not. Specifically, obligation models distinguish four possible states for duties: invalid (when at the point a duty is issued the completion point has already passed); fulfilled (when the obligation is carried out within the associated interval); violated (when it is not carried out within the associated temporal interval) and pending (when the obligation has not yet been carried, but the interval is still valid). In some cases, \(p\)'s duty can be considered fulfilled even if \(p\) has not performed the required action, because another principal has done it. This is the case in the example above, where all members of the security team have the obligation to call the fire department before deactivating the alarm, but the obligation is fulfilled as soon as one of them makes the call. In order to distinguish both kinds of obligations, we will call **individual obligations** those that have to be fulfilled necessarily by the principal to whom the obligation is assigned, and **collective obligations** those where the obligation is assigned to several principals and can be fulfilled by any of the principals in the group.

### 4 Obligations in the category-based metamodel

The notion of a category will be used to specify obligations that apply to groups of principals. However, the groupings of principals for authorisations and for obligations are not necessarily the same (for instance, the category UG student is assigned a certain set of permissions, which all UG students enjoy, whereas some UG students belong to the home student category and others to the international student category, with different obligations). We call permission categories those used to specify authorisations, and obligation categories those used to specify duties. To model obligations and duties we extend the metamodel
to include the following sets of entities and relations in addition to the ones defined in section 2:

- Countable sets $\mathcal{E}$ and $\mathcal{GE}$ of events and event types, denoted by $e, e_0, e_1, \ldots$ and $ge, ge_0, ge_1, \ldots$, respectively.
- A countable set $\mathcal{H}$ of event histories, denoted by $h, h_0, h_1, \ldots$.
- Obligation-category assignment: $\mathcal{OARCA} \subseteq \mathcal{A} \times \mathcal{R} \times \mathcal{GE} \times \mathcal{GE} \times \mathcal{C}$, such that $(a, r, ge_1, ge_2, c) \in \mathcal{OARCA}$ if the action $a \in \mathcal{A}$ on resource $r \in \mathcal{R}$ must be performed in the interval defined by two events of type $ge_1, ge_2$ by principals assigned to the category $c \in \mathcal{C}$. To accommodate individual and collective obligation assignments, this relation is partitioned into two relations: $\mathcal{OARCA}_I$ and $\mathcal{OARCA}_C$. Thus, $(a, r, ge_1, ge_2, c) \in \mathcal{OARCA}_I$ if every member of the category $c$ must perform the action $a \in \mathcal{A}$ on resource $r \in \mathcal{R}$ in the interval defined by $ge_1, ge_2$, and $(a, r, ge_1, ge_2, c) \in \mathcal{OARCA}_C$ if any member of $c$ must perform the action $a \in \mathcal{A}$ on resource $r \in \mathcal{R}$ in the interval defined by $ge_1, ge_2$ and it is sufficient that one of them does it.
- Obligation-principal assignment: $\mathcal{OPAR} \subseteq \mathcal{P} \times \mathcal{A} \times \mathcal{R} \times \mathcal{GE} \times \mathcal{GE}$, such that $(p, a, r, ge_1, ge_2) \in \mathcal{OPAR}$ if a principal $p \in \mathcal{P}$ must perform the action $a \in \mathcal{A}$ on the resource $r \in \mathcal{R}$ in the interval defined by two events of type $ge_1, ge_2$. If individual and collective obligations are modelled, then this relation is replaced by two relations: $\mathcal{OPAR}_I \subseteq \mathcal{P} \times \mathcal{A} \times \mathcal{R} \times \mathcal{GE} \times \mathcal{GE}$, defining individual obligations, and $\mathcal{OPAR}_C \subseteq \mathcal{P} \times \mathcal{A} \times \mathcal{R} \times \mathcal{GE} \times \mathcal{GE} \times \mathcal{C}$ for collective obligations, such that $(p, a, r, ge_1, ge_2) \in \mathcal{OPAR}_I$ iff principal $p \in \mathcal{P}$ must perform the action $a \in \mathcal{A}$ on the resource $r \in \mathcal{R}$ in the interval defined by two events of type $ge_1, ge_2$, and $(p, a, r, ge_1, ge_2, c) \in \mathcal{OPAR}_C$ if principal $p \in \mathcal{P}$ or any other member of $c$ must perform the action $a \in \mathcal{A}$ on the resource $r \in \mathcal{R}$ in the interval defined by two events of type $ge_1, ge_2$.
- Duty: $\mathcal{DPAR} \subseteq \mathcal{P} \times \mathcal{A} \times \mathcal{R} \times \mathcal{E} \times \mathcal{E} \times \mathcal{H}$, such that $(p, a, r, e_1, e_2, h) \in \mathcal{DPAR}$ if a principal $p \in \mathcal{P}$ must perform the action $a \in \mathcal{A}$ on the resource $r \in \mathcal{R}$ in the interval between the events $e_1$ and $e_2$ in the event history $h$. To accommodate individual and collective obligations, this relation is partitioned into $\mathcal{DPAR}_I$ and $\mathcal{DPAR}_C$, similarly to $\mathcal{OPAR}$.
- Event Typing: $\mathcal{ET} \subseteq \mathcal{E} \times \mathcal{GE} \times \mathcal{H}$, such that $(e, ge, h) \in \mathcal{ET}$ if the event $e$ is an instance of $ge$ in $h$. This will be abbreviated as $h \vdash e : ge$.
- Event Interval: $\mathcal{EI} \subseteq \mathcal{E} \times \mathcal{E} \times \mathcal{H}$, such that $(e_1, e_2, h) \in \mathcal{EI}$ if the event $e_2$ closes the interval started by the event $e_1$ in $h$.

### 4.1 Obligation Axioms

In the metamodel, the relations $\mathcal{OPAR}$ and $\mathcal{DPAR}$ are derivable from the others. They satisfy the following core axioms, where we assume that there exists
a relationship $\subseteq_o$ between obligation categories; this can simply be equality, set inclusion (the set of principals assigned to $c \in C$ is a subset of the set of principals assigned to $c' \in C$), or an application specific relation may be used.

\[
(o1) \forall p \in P, \forall a \in A, \forall r \in R, \forall ge_1, ge_2 \in GE, \forall c, c' \in C, (p, c) \in PCA \land c \subseteq_o c' \land (a, r, ge_1, ge_2, c') \in OARCA \iff (p, a, r, ge_1, ge_2) \in OPAR
\]

\[
(o2) \forall p \in P, \forall a \in A, \forall r \in R, \forall e_1, e_2 \in E, \forall h \in H, \forall c \in C, (p, a, r, ge_1, ge_2) \in OPAR, \forall c \in C, (p, c) \in PCA \land c \subseteq_o c' \land (a, r, ge_1, ge_2, c') \in OARCA \iff (p, a, r, e_1, e_2, h) \in DPAR
\]

The axiom $o1$ is the essence of the category-based metamodel: it specifies that the principals that are members of a category to which the obligation $(a, r, ge_1, ge_2)$ has been assigned have this obligation. The axiom $o2$ shows how to derive duties. The relation $\subseteq_o$ specifies a hierarchy between obligation categories, which does not necessarily correspond to the way permissions are inherited (specified by the relation $\subseteq$ in axiom $(a1)$).

To accommodate collective and individual obligations, the following variants of the axioms should be included.

\[
(o1_I) \forall p \in P, \forall a \in A, \forall r \in R, \forall ge_1, ge_2 \in GE, \forall c, c' \in C, (p, c) \in PCA \land c \subseteq_o c' \land (a, r, ge_1, ge_2, c') \in OARCA_I \iff (p, a, r, ge_1, ge_2) \in OPAR_I
\]

\[
(o1_C) \forall p \in P, \forall a \in A, \forall r \in R, \forall ge_1, ge_2 \in GE, \forall c \in C, (p, c) \in PCA \land c \subseteq_o c' \land (a, r, ge_1, ge_2, c') \in OARCA_C \iff (p, a, r, ge_1, ge_2, c') \in OPAR_C
\]

\[
(o1') OPAR_I \cap OPAR_C = \emptyset
\]

where $OPAR_C$ represents the projection of the relation which discards the last column (the category). The variants of the axiom $o2$ are defined similarly. Axiom $(o1')$ indicates that the same obligation cannot be both collectively and individually fulfilled. Additionally, the relation $OARCA_C$ should not assign the same obligation to two categories in the $\subseteq_o$ relation, to avoid generated redundant tuples in $OPAR_C$.

### 4.2 Event and event type representation

To consider examples of obligation policies we will describe a possible representation for events and event types. Following the event calculus approach [12], we consider events as action occurrences, or action happenings. Actions
can be either elementary or compound, i.e. consisting of sets of (simultaneously happened) elementary actions [9]. We use letters $a$ and $e$, possibly with indexes, to identify actions and events, respectively. We will use a binary representation for events, introduced in [4] and partly motivated by Davidson’s work on event semantics [8]. This choice provides a most flexible representation with negligible computational overheads. In our case the set of arguments depends, not only on the action involved but also on the context where the event description is used.

**Definition 5 (Event).** An event description is a finite set of ground 2-place facts (atoms) that describe an event, uniquely identified by $e_i, i \in \mathbb{N}$, and which includes two necessary facts, written $\text{happens}(e_i, t_j)$ and $\text{act}(e_i, a_l)$, and $n$ non-necessary facts ($n \geq 0$).

In [4] an event description includes three necessary facts, and $n$ non-necessary facts ($n \geq 0$). Unlike [4], we use only two types of necessary facts. Their intended meanings may be described as follows: $\text{happens}(e_i, t_j)$ means the event $e_i$ happens at time $t_j$; $\text{act}(e_i, a_l)$ means the event $e_i$ involves an action $a_l$. The events should be positioned in the history in the order in which they happened.

Sometimes we will consider the only predicate $\text{act}$ as necessary and $\text{happens}$ as optional. In this case, the history would include both events with specified time and events without indicating time. Whether we use the $\text{happens}$ predicate or not, the history should reflect the order of the events.

If the event $e_i$ involves the subject $s_m$, then the corresponding non-necessary fact $\text{subject}(e_i, s_m)$ can be added to the event description when we need this fact. In [4], $\text{subject}(e_i, s_m)$, would be the third necessary fact. Similarly, if the event $e_i$ involves the resource $r_l$, then the fact $\text{object}(e_i, r_l)$ can be added to the description. And so on. Thus, the event description given by the set \{ $\text{happens}(e_i, t_j)$, $\text{act}(e_i, a_l)$, $\text{subject}(e_i, s_m)$, $\text{object}(e_i, r_l)$ \} represents a happening $e_i$ at a time $t_j$ of an act $a_l$ performed by a subject $s_m$ on a resource $r_l$.

To define obligations associated with events, here we distinguish between event types, denoted by $ge_i$ and specific events $e_i$. As it was noted earlier, in real modelling the set of non-necessary predicates involved in a description of an event is determined not only by the action assigned to this event but also by current context. We define an event type, or generic event, to consist of the necessary fact indicating an action and a set of predicates over expressions with variables including another necessary predicate $\text{happens/2}$.

**Definition 6 (Event type and Instance).** A set $ge$ of binary atoms represents an event type, also called generic event, if there exists a substitution $\sigma$ such that the instantiation $ge\sigma$ is an event description (see Definition 5).

---

1 We restrict attention to events without a duration.
An event \( e \) is an instance of a generic event \( ge \), denoted as \( e : ge \), if there is a substitution \( \sigma \) such that \( ge\sigma \subseteq e \). In other words, if \( ge \) subsumes \( e \).

The action of a substitution on a generic event may simply be a syntactic operation, or, depending on the application and the kind of data used to define events, instantiation may require some computation (we call it a semantic instantiation in the latter case). We give examples below.

**Example 2.** The event \( \{ \text{happens}(e_2, 12.25), \text{act}(e_2, \text{activate}), \text{object}(e_2, \text{alarm}) \} \) is an instance of the generic event

\[
\text{alarmON} = \{ \text{happens}(E_2, T), \text{act}(E_2, \text{activate}), \text{object}(E_2, \text{alarm}) \}
\]

with the substitution \( \sigma = \{ E_2 \mapsto e_2, T \mapsto 12.25 \} \).

The event \( \{ \text{happens}(e_3, 12.45), \text{act}(e_3, \text{deactivate}), \text{object}(e_3, \text{alarm}) \} \) is an instance of the generic event

\[
\text{alarmOFF} = \{ \text{happens}(E, T + D), \text{act}(E, \text{deactivate}), \text{object}(E, \text{alarm}) \}
\]

with the substitution \( \sigma = \{ E \mapsto e_3, T \mapsto 12.25, D \mapsto 20 \} \). Here we are using a semantic instantiation function.

**Example 3.** Assume that in a given university, every international student (i.e., a student who comes from a different country) must have a valid visa before starting the course. We define a category “international-student”, such that

\[
(\text{obtainVisa}, \text{passport}, \bot, \text{course-registration}, \text{international-student}) \in OARCA^2.
\]

Here the event type that initiates the obligation is not specified (\( \bot \)), but the final one is (\( \text{course-registration} \)).

**Example 4.** Consider a hospital, where doctors are allowed to read and write the medical records of the patients in their department. Assume that the policy in place in the hospital also states that any doctor is allowed to read the medical record of a patient who is in an emergency situation (even if the patient is not in their department) but in that case they or any doctor in their department must declare this access (by creating an entry in the hospital log). Let \( \text{patient} \) be a category consisting of all patients (of a given hospital), and \( \text{doctor} \) be a category consisting of all doctors (of the given hospital). Let \( \text{patient}(X) \) be a (parameterised) category consisting of all patients in the department \( X \), and \( \text{doctor}(X) \) be a (parameterised) category consisting of all doctors in the department \( X \), such that for all \( X \), \( \text{doctor}(X) \subseteq \text{doctor} \), i.e., the category \( \text{doctor}(X) \) inherits all permissions from the category \( \text{doctor} \) and similarly \( \text{patient}(X) \subseteq \text{patient} \). Assume

\[2\] or \( OARCA_t \) if we need to distinguish between individual and collective obligations.
that these categories and hierarchical relation are used both for permissions and for obligations.

To model this scenario, we assume the relations $\mathcal{PCA}$ and $\mathcal{ARCA}$ satisfy the following axioms, where $\text{emerg}(\text{bcrd}, P)$ is true if an event $\text{bcrd}$ initiating a cardiac emergency for $P$ has been detected, and no event ending the emergency has been recorded:

\[ \forall P, (P, \text{patient}(D)) \in \mathcal{PCA} \Rightarrow (\text{read, record}(P), \text{doctor}(D)) \in \mathcal{ARCA} \]

\[ \forall P, (P, \text{patient}(D)) \in \mathcal{PCA} \land \text{emerg}(\text{bcrd}, P) \Rightarrow (\text{read, record}(P), \text{doctor}(D)) \in \mathcal{ARCA} \]

Moreover, we include the following axiom for $\mathcal{OARCA}_C$, where $\text{gen\_read}(X, P)$ is a generic event representing the act of $\text{doctor}(X)$ reading the medical record of $\text{patient}(P)$:

\[ \forall P, \forall X, \forall D, \forall D' ((P, \text{patient}(D)) \in \mathcal{PCA}, (X, \text{doctor}(D')) \in \mathcal{PCA}, D \neq D') \Rightarrow (\text{declare, admin-log, gen\_read}(X, P), \bot, \text{doctor}(D')) \in \mathcal{OARCA}_C) \]

5 A Rewrite Semantics for Obligations

Operationally, the axioms $o1$ and $o2$ given in Section 4 can be realised through a set of function definitions. In this section we assume all the obligations are individual; the extension to accommodate individual and collective obligations is straightforward. The information contained in the relations $\mathcal{PCA}$ and $\mathcal{ARCA}$ is modelled by the functions $\text{pca}$ and $\text{oarca}$, respectively, where $\text{pca}$ returns the list of all the categories assigned to a principal (for efficiency reasons, a separate function $\text{opca}$ could be defined to compute obligation categories) and $\text{oarca}$ returns the list of obligations assigned to a category, e.g., defined by the rule:

\[
\text{oarca}(c) \to [(a_1, r_1, ge_{1_1}, ge_{1_2}), \ldots, (a_n, r_n, ge_{n_1}, ge_{n_2})].
\]

We assume that the function $\text{oarca}^*$ takes as input a list of obligation categories and computes the list of obligations for all the categories in the list (similar to $\text{arca}^*$, see Section 2).

In addition, we assume that the functions $\text{type}$ and $\text{interval}$, specific to the particular system modelled, are also available. The function $\text{type}$ implements the typing relation for events, that is, it computes the event type for a given event $e$ in $h$ (taking into account the semantic instantiation relation associated with the events of interest). The function $\text{interval}$ implements the relation $EI$, which links an event with the event that closes the interval in $h$.

**Definition 7.** The rewrite-based specification of the axiom $(o1)$ in section 4.1 is given by the rewrite rule:

\[
(1) \text{opar}(p, a, r, ge_1, ge_2) \to (a, r, ge_1, ge_2) \in \text{oarca}^*(\text{obelow}(\text{opca}(p)))
\]
where the function $\in$ is a membership operator on lists, and, as the function name suggests, obelow computes the set of categories that are below (w.r.t. the hierarchy defined by the $\subseteq_o$ relation) any of the categories given in the list $\text{opca}(p)$. For example, for a given category $c$, this could be achieved by using a rewrite rule $\text{obelow}([c]) \rightarrow [c, c_1, \ldots, c_n]$. Intuitively, this means that if $c'$ is below $c$, then $c$ inherits all the obligations of $c'$.

The rewrite-based specification of the axiom $(o2)$ is given by:

$$(r2) \quad \text{duty}(p, a, r, e_1, e_2, h) \rightarrow \text{opar}(p, a, r, \text{type}(e_1, h), \text{type}(e_2, h)) \quad \text{and}$$

$$\text{interval}(e_1, e_2, h)$$

where the functions $\text{type}$ and $\text{interval}$ are specific to the system modelled, as mentioned above. Additionally, we consider the following function to check obligations:

$$(r3) \quad \text{eval-obligation}(p, a, r, ge_1, ge_2, h) \rightarrow \begin{cases} \text{if } \text{opar}(p, a, r, ge_1, ge_2) \text{ then} \\ \text{append}(\text{chk-closed}^*(\text{closed}(ge_1, ge_2, h), p, a, r), \text{chk-open}^*(\text{open}(ge_1, ge_2, h), p, a, r)) \\ \text{else } \text{not-applicable} \end{cases}$$

where the function $\text{append}$ is a standard function that concatenates two lists, $\text{closed}$ computes the sublists of $h$ that start with an event $e_1$ of type $ge_1$ and finish with an event $e_2$ of type $ge_2$ that closes the interval for this obligation, and similarly $\text{open}$ returns the subhistories of $h$ that start with the event $e_1$ of type $ge_1$ and for which there is no event $e_2$ of type $ge_2$ in $h$ that closes the interval for this obligation.

The function $\text{chk-closed}$ with inputs $h', p, a, r$ checks whether in the subhistory $h'$ there is an event where the principal $p$ has performed the action $a$ on the resource $r$, returning a result $\text{fulfilled}$ if that is the case, and $\text{violated}$ otherwise.

The function $\text{chk-open}$ with inputs $h', p, a, r$ checks whether in the subhistory $h'$ there is an event where the principal $p$ has performed the action $a$ on the resource $r$, returning a result $\text{fulfilled}$ if that is the case, and $\text{pending}$ otherwise.

The functions $\text{chk-closed}^*$ and $\text{chk-open}^*$ do the same but for each element of a list of subhistories, returning a list of results.

According to the previous specification, if the obligation $(a, r, ge_1, ge_2)$ does not concern $p$ then $\text{eval-obligation}(p, a, r, ge_1, ge_2, h)$ returns $\text{not-applicable}$, otherwise, the function $\text{eval-obligation}$ returns a list of results containing the status of $p$ in relation to this obligation according to $h$. Usually, $h$ will be event history that is relevant to the scenario being modelled. For instance, it could be the full history of events in the system, or we could restrict ourselves to the events in the last year, or the events that happened during a specific interval. In our model, it is not possible for a duty to be invalid (thanks to axiom $(o2)$).
If the system modelled includes collective and individual obligations, then the functions \texttt{chk-closed} and \texttt{chk-open} should take into account this in order to check if the obligation has been fulfilled.

\textit{Example 5.} Considering again the hospital scenario mentioned in Example 4. Assume an investigation is taking place to establish if Doctor Tom Smith, who treated patient John Lewis in an emergency situation occurring in November 2012, but is not John Lewis’s doctor, has fulfilled his obligation with respect to declaring the access to this patient’s records. This is a collective obligation which can be discharged by any of the doctors in his department. Accordingly, we assume the functions \texttt{chk-closed} and \texttt{chk-open} perform collective checks.

In this case, we simply evaluate the term

\[
\texttt{eval-obligation}(\text{TomSmith, declare, admin-log, ge, } \perp, h)
\]

where \(ge = \{\text{act}(E, \text{read-pr}(\text{JohnLewis}), \text{sub}(E, \text{TomSmith})\} \) and \(h\) is the event history that starts on the 1st November 2012 and ends today.

\begin{subsection}{5.1 Analysis of policies}

Specifying policies via term rewriting systems, which have a formal semantics, has the advantage that this representation admits the possibility of proving properties of policies. In previous work, rewriting techniques were used to derive properties of authorisation policies in the metamodel (see, e.g., [5]). In particular, rewriting properties such as confluence (which implies the unicity of normal forms) and termination (which implies the existence of normal forms for all terms) were used to prove totality and consistency of authorisation policies. Here we apply similar techniques to prove properties of policies that include obligations.

Given a policy specification, \texttt{eval-obligation}(\(p, a, r, ge_1, ge_2, h\)) should return an answer, either indicating that the obligation \((a, r, ge_1, ge_2)\) does not apply to \(p\), or indicating whether this obligation was fulfilled, violated or is still pending according to \(h\). Of course, if \(h\) is sufficiently long and depending on the events it contains, it is possible that at different points the obligation was fulfilled, violated or pending. This is why the function \texttt{eval-obligation}(\(p, a, r, ge_1, ge_2, h\)) returns a list of results.

Termination and confluence of the rewrite system specifying \texttt{opar}, \texttt{duty} and \texttt{eval-obligation} is therefore a minimum requirement in the metamodel. These properties will depend on the specific functions defined for the specific policy modelled, and unfortunately they are generally undecidable. However, sufficient conditions exist and tools such as CiME [7] could be used to help checking properties of the first-order rewrite rules.
The metamodel is sufficiently expressive to permit the definition of policies where authorisations and obligations co-exist and may depend on each other. In this case, security administrators need to check the consistency of the definitions: the policy should ensure that every user has the required permissions in order to fulfill all his/her duties. In practice, the authorisation relation $\mathcal{P \cdot AR}$ may be defined using categories whose definition takes into account dynamic conditions (e.g., $h$). In order to ensure that only duties that are consistent with authorisations are issued, the operational semantics could be modified, by including in the right-hand side of the rule $r2$ the condition $\text{par}(p, a, r, h, e_1, e_2)$, where the function $\text{par}$ computes the authorisation relation within the interval defined by $e_1$ and $e_2$ in $h$.

\section{6 Related work}

Several models dealing with the notion of obligations have been proposed in the literature [16, 10, 14, 13, 11, 6] (amongst others), with different features and providing different levels of abstraction in the definition and management of obligations. Some systems consider obligations to be fulfilled by the systems and therefore never violated, whereas others consider obligations to be fulfilled by users, in which case other questions arise such as how can the system guarantee that the user will be allowed to fulfill the obligations, and the notion of accountability (if a user, with the necessary permissions does not fulfill its obligations then he becomes accountable).

The framework in [13] deals both with system and user obligations and measures to enforce the fulfillment of obligations. Obligations are triggered by time or events and violated obligations are treated by the system.

The system presented in [11] extends the usage control (UCON) model with obligations and separates system obligations (which are called trusted obligations) from user obligations (called non-trusted). Mechanisms are proposed to deal with the non-trusted obligations. The system does not consider the interaction of obligations with permissions, neither deals with conditions on obligations. In [15], the ABC core model was defined for UCON, dealing with authorizations, obligations and conditions. This approach differs from ours, since it describes not a general metamodel, but instead a family of models depending on the type of authorizations, obligations and conditions that are considered and the necessary components for each model in the family. Our approach also considers authorizations, obligations and conditions, but in an uniform way.

The approach followed in [6] was to formalise policies (using Horn clauses) trying to minimise the number of actions that need to be performed before allowing an access request. Although the initial system only dealt with sys-
tems without violated obligations, this was later extended to handle violation of obligations. This approach was limited in the type of obligations that could be modelled, because it lacked appropriate mechanisms to deal with temporal constraints.

In [10] a model is presented for authorisation systems dealing with the notion of accountability; obligations are defined as an action on a resource by a user within a time interval (defined in clock ticks from a predetermined starting point). The monitoring and restoring of accountability was further explored in [16], where a transition relation is defined and conditions to guarantee accountability are established. The notion of obligations defined in these works corresponds to concrete obligations (duties) in our model, and although this is not the focus of this paper, we believe that rewriting can be used to verify the properties of accountability studied in these papers.

A more general model dealing with obligations and its relation to access control and privacy policies was defined in [14]. The model also investigates the interaction of obligations and permissions and its undesirable effects such as obligation cascading. The category-based model is expressive enough to model the concepts of obligations defined in this work, and rewriting properties can be used to further explore the interplay between permissions and obligations.

Note also that the CBAC model we consider here is an extension of the model defined in [5], therefore all the results regarding the rewriting semantics are also valid for this model, but are not the focus of this paper.

7 Conclusions

We augmented the CBAC metamodel with a general notion of obligation. The framework is expressive enough to deal with most of the features relevant to authorisations and obligations and provides means to reason about them.

Our model distinguishes individual and collective obligations by partitioning the set of obligations into two distinct sets. However, there are some situations where it can be difficult to distinguish between one or the other type. An alternative approach to be investigated in the future is to base this distinction on the definition of the obligation-category assignment relation. That is, although at a certain moment all the principals belonging to a particular category are obliged to perform some action, the definition of the category can depend on the fact of the action not having been performed at the time. In future work we will formally define the typing relation for events. We also wish to explore appropriate mechanisms to deal with accountability.
Access control and obligations in the category-based metamodel

References

Abstract. Symbolic execution extends concrete execution by allowing symbolic input data and then exploring all feasible execution paths. It has been defined and used in the context of many different programming languages and paradigms. A symbolic execution engine is at the heart of many program analysis and transformation techniques, like partial evaluation, test case generation or model checking, to name a few. Despite its relevance, traditional symbolic execution also suffers from several drawbacks. For instance, the search space is usually huge (often infinite) even for the simplest programs. Also, symbolic execution generally computes an overapproximation of the concrete execution space, so that false positives may occur. In this paper, we propose the use of a variant of symbolic execution, called concolic execution, for test case generation in Prolog. We argue that this technique computes an underapproximation of the concrete execution space (thus avoiding false positives) and scales up better to medium and large Prolog applications.

1 Introduction

There is a renewed interest in symbolic execution [8, 2], a well-known technique for program verification, testing, debugging, etc. In contrast to concrete execution, symbolic execution considers that the values of some input data are unknown, i.e., some input parameters $x,y,...$ take symbolic values $X,Y,...$. Because of this, symbolic execution is often non-deterministic: at some control statements, we need to follow more than one execution path because the available information does not suffice to determine the validity of a control expression, e.g., symbolic execution may follow both branches of the conditional “if $(x>0)$ then $exp1$ else $exp2$” when the symbolic value $X$ of variable $x$ is not constrained enough to imply neither $x>0$ nor $\neg(x>0)$. Symbolic states include a path condition that stores the current constraints on symbolic values, i.e., the

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conditions that must hold to reach a particular execution state. E.g., after symbolically executing the above conditional, the derived states for \( \text{exp1} \) and \( \text{exp2} \) would add the conditions \( X > 0 \) and \( X \leq 0 \), respectively, to their path conditions.

Traditionally, formal techniques based on symbolic execution have enforced *soundness*: if a symbolic state is reached and its path condition is satisfiable, there must be a concrete execution path that reaches the corresponding concrete state. In contrast, we say that symbolic execution is *complete* when every reachable state in a concrete execution is “covered” by some symbolic state. For the general case of infinite state systems, completeness usually requires some kind of *abstraction* (as in infinite state model checking).

In the context of logic programming, we can find many techniques that use some form of *complete* symbolic execution, like partial evaluation [11, 12, 3]. However, these overapproximations of the concrete semantics may have a number of drawbacks in the context of testing and debugging. On the one hand, one should define complex subsumption and abstraction operators since the symbolic search space is usually infinite. These abstraction operators, may introduce *false positives*, which is often not acceptable when debugging large applications. On the other hand, because of the complexity of these operators, the associated methods usually do not scale to medium and large applications.

In imperative programming, an alternative approach, called *concolic execution* [5, 15], has become popular in the last years. Basically, concolic execution proceeds as follows: first, a concrete execution using random input data is performed. In parallel to the concrete execution, a symbolic execution is also performed, but restricted to the same conditional choices of the concrete execution. Then, by negating one of the constraints in the symbolic execution, new input data are obtained, and the process starts again. Here, only concrete executions are considered and, thus, no false positives are produced. This approach has given rise to a number of powerful and scalable tools in the context of imperative and concurrent programming, like Java Pathfinder [13] and SAGE [6].

In this paper, we present a novel scheme for testing pure Prolog (without negation) based on a notion of concolic execution. To the best of our knowledge, this is the first approach to concolic execution in the context of a declarative programming paradigm.

2 Preliminaries

We assume some familiarity with the standard definitions and notations for logic programs [10]. Nevertheless, in order to make the paper as self-contained as possible, we present in this section the main concepts which are needed to understand our development.

In this work, we consider a first-order language with a fixed vocabulary of predicate symbols, function symbols, and variables denoted by \( \Pi, \Sigma \) and \( \mathcal{V} \), respectively. In the following, we let \( \overline{o} \) denote the sequence of syntactic objects \( o_1, \ldots, o_n \). Also, we often use \( \overline{\sigma} \) when the number of elements in the sequence is irrelevant. We let \( \mathcal{T}(\Sigma, \mathcal{V}) \) denote the set of terms constructed using symbols
from $\Sigma$ and variables from $V$. An atom has the form $p(t_1) \ldots p(t_n)$ with $p(n) \in \Pi$ and $t_i \in T(\Sigma, V)$ for $i = 1, \ldots, n$. A goal is a finite sequence of atoms $A_1, \ldots, A_n$, where the empty goal is denoted by true. A clause has the form $H \rightarrow B$ where $H$ is an atom and $B$ is a goal. A logic program is a finite sequence of clauses. A syntactic object $s$ is ground if $\text{Var}(s) = \emptyset$. In this work, we only consider finite ground terms.

Substitutions and their operations are defined as usual. In particular, the set $\text{Dom}(\sigma) = \{ x \in V \mid \sigma(x) \neq x \}$ is called the domain of a substitution $\sigma$. We let $id$ denote the empty substitution. The application of a substitution $\theta$ to a syntactic object $s$ is usually denoted by juxtaposition, i.e., we write $s\theta$ rather than $\theta(s)$. A syntactic object $s_1$ is more general than a syntactic object $s_2$, denoted $s_1 \preceq s_2$, if there exists a substitution $\theta$ such that $s_2 = s_1\theta$. A variable renaming is a substitution that is a bijection on $V$. Two syntactic objects $t_1$ and $t_2$ are variants (or equal up to variable renaming), denoted $t_1 \equiv t_2$, if $t_1 = t_2\rho$ for some variable renaming $\rho$. A substitution $\theta$ is a unifier of two syntactic objects $t_1$ and $t_2$ if $t_1\theta = t_2\theta$; furthermore, $\theta$ is the most general unifier of $t_1$ and $t_2$, denoted by $\text{mgu}(t_1, t_2)$ if, for every other unifier $\sigma$ of $t_1$ and $t_2$, we have that $\theta \preceq \sigma$.

The notion of computation rule $R$ is used to select an atom within a goal for its evaluation. Given a program $P$, a goal $G \equiv A_1, \ldots, A_n$, and a computation rule $R$, we say that $G \sim_{P,R,\sigma} G'$ is an SLD resolution step for $G$ with $P$ and $R$ if

- $R(G) = A_i$, $1 \leq i \leq n$, is the selected atom,
- $H \rightarrow B$ is a renamed apart clause of $P$ (in symbols $H \rightarrow B \ll P$),
- $\sigma = \text{mgu}(A,H)$, and
- $G' \equiv (A_1, \ldots, A_{i-1}, B, A_{i+1}, \ldots, A_n)\sigma$.

We often omit $P$, $R$ and/or $\sigma$ in the notation of an SLD resolution step when they are clear from the context. An SLD derivation is a (finite or infinite) sequence of SLD resolution steps. We often use $G_0 \sim_{\theta_0}^* G_n$ as a shorthand for $G_0 \sim_{\theta_0} G_1 \sim_{\theta_1} \ldots \sim_{\theta_n} G_n$ with $\theta = \theta_n \circ \ldots \circ \theta_1$ (where $\theta = \{ \}$ if $n = 0$). An SLD derivation $G \sim_{\theta}^* G'$ is successful when $G' = \text{true}$; in this case, we say that $\theta$ is the computed answer substitution. SLD derivations are represented by a (possibly infinite) finitely branching tree.

3 A Deterministic Semantics

In this section, we introduce a deterministic small-step semantics for pure Prolog (without negation). Basically, as we will see in the next section, we need to keep some information through the complete Prolog computation, and the usual semantics based on non-determinism and backtracking is not adequate for this purpose. Therefore, we propose the use of a stack to store alternative execution paths that are tried when a failure is reached. The resulting small-step semantics is clearly equivalent to the original one when a depth-first search is considered.
Actually, our deterministic semantics is essentially equivalent to (a subset of) the linear semantics presented in [16].

In the following, we assume that the program clauses are labeled. In particular, given a program \( P \), we use the notation \( H \leftarrow B \ll P \) to refer to a (renamed apart) labeled clause \( H \leftarrow B \) in \( P \). Labels must be unique. Moreover, we only consider Prolog’s left-to-right computation rule, and assume that only the computation of the first answer is relevant (as it is common in practical Prolog applications).

Our semantics deals with states, which are defined as follows:

\textbf{Definition 1 (state).} A state is a pair \( \langle G; S \rangle \), where \( G \) is a goal and \( S \), the stack, is a (possibly empty) list of tuples \((\ell; \sigma; G')\) with \( \ell \) a clause label, \( \sigma \) a substitution, and \( G' \) a goal.

The small-step deterministic semantics is defined as the smallest relation that obeys the labeled transition rules shown in Figure 1, where \([H|R]\) denotes a list with head \( H \) and tail \( R \), and “++” denotes list concatenation.

Given a goal \( G_0 \), the initial state has the form \( \langle G_0; [] \rangle \). The transition relation is labeled with the substitution computed in the step and either \( u(\ell_n) \), denoting an unfolding step with the clauses labeled with \( \ell_n \), \( b(\ell) \), denoting a backtracking step that tries a clause labeled with \( \ell \), or \( f \), denoting a failing derivation.

Let us briefly explain the rules of the small-step semantics:

- The unfolding rule proceeds as in standard SLD resolution, but considers all matching clauses, so that all SLD resolution steps are performed in one go. The first unfolding step is used to replace the goal component of the state, while the remaining ones (if any) are added on top of the stack (thus we mimick the usual depth-first search of Prolog). Here, the labels of the clauses and the computed unifiers are also stored in the stack in order to recover this information when a backtracking step is performed.
The backtracking rule applies when no further unfolding is possible and the goal component is not true (the empty goal). In this case, we discard the current goal and consider the first goal in the stack, using the clause label and the substitution to label the transition step.

Finally, the failure rule is used to terminate a computation that reaches a goal in which the selected atom does not match any rule and, moreover, there are no alternatives in the stack.

A successful computation has the form \( \langle G_0; \emptyset \rangle \xrightarrow{\sigma_1} \langle G_1; S_1 \rangle \xrightarrow{\sigma_2} \ldots \xrightarrow{\sigma_n} \langle \text{true}; S_n \rangle \), where \( \sigma_1, \sigma_2, \ldots, \sigma_n \) is the computed answer substitution. A failing computation has the form \( \langle G_0; \emptyset \rangle \xrightarrow{\sigma_1} \langle G_1; S_1 \rangle \xrightarrow{\sigma_2} \ldots \xrightarrow{\sigma_n} \langle G_n; S_n \rangle \xrightarrow{id} \langle \text{fail}; \emptyset \rangle \).

Now, we introduce the following notion of execution trace, that will be used in the next section to steer the symbolic execution.

Definition 2 (trace). Let \( \langle G_0; S_0 \rangle \xrightarrow{\sigma_1} \langle G_1; S_1 \rangle \xrightarrow{\sigma_2} \ldots \xrightarrow{\sigma_n} \langle G_n; S_n \rangle \) be a computation. The associated trace is the list \( [s_1, s_2, \ldots, s_n] \), where each \( s_i \) is either of the form \( u(\ell m) \), \( b(\ell) \) or \( f \).

Example 1. Consider the rev_acc_type program to reverse a list using an accumulator and also checking the type of the input parameter (from the DPPD library [9]), extended with predicates main, length, and foo:

\[
\begin{align*}
(1) & \quad \text{main}(L,N,R) :- \\
& \quad \text{length}(L,N), \\
& \quad \text{rev}(L,\emptyset,R), \\
& \quad \text{foo}(a).
(2) & \quad \text{main}(L,_,N,\text{error}).
(3) & \quad \text{rev}(\emptyset,A,A).
(4) & \quad \text{rev}(H\left| T\right.,\text{Acc},\text{Res}) :- \\
& \quad \text{is_list}(\text{Acc}), \\
& \quad \text{rev}(T,[H|\text{Acc}],\text{Res}).
(5) & \quad \text{is_list}([\]).
(6) & \quad \text{is_list}([\_H|T]) :- \\
& \quad \text{is_list}(T).
(7) & \quad \text{length}([],0).
(8) & \quad \text{length}([\_H\_R],s(N)) :- \\
& \quad \text{length}(R,N).
(9) & \quad \text{foo}(b).
\end{align*}
\]

Here, we use natural numbers as clause labels. Predicate main considers two cases: if the input list \( L \) has length \( N \) (the length is represented using natural numbers built from 0 and \( s(\_) \) to avoid the use of built-ins), the reverse of \( L \) is computed; otherwise, we assume that an error occurs. The computation for the initial goal \( \text{main}([a,b],s(s(0)),R) \) is shown in Figure 2, where only the relevant computed substitutions are shown. The trace associated to the computation is 

\( [u(1,2), u(8), u(8), u(7), u(4), u(5), u(4), u(6), u(5), u(3), b(2)] \)

4 Concolic Execution

In this section, we introduce the semantics of concolic execution. Essentially, it deals with symbolic input data (free variables in our context), as in standard
symbolic execution, but is driven by a concrete execution. Often, a single algorithm mixing both concrete and symbolic execution is introduced. In contrast, for clarity, we prefer to keep both calculi independent: the concrete semantics produces a trace, which is then used to steer the symbolic execution.¹

The symbolic states for concolic execution are defined as follows:

**Definition 3 (symbolic state).** A symbolic state is a tuple \(\langle \tau; L; G; S; T \rangle\), where

- \(\tau\) is a computation trace,
- \(L\) is a list of clause labels (namely, a stack that keeps track of the current clause environment),²
- \(G\) is a goal,
- \(S\) is a (possibly empty) list of tuples \(\langle \ell; L'; \sigma; G' \rangle\), where \(L'\) is also a list of clause labels, and
- \(T\) is a set of clause labels (the labels of those clauses not yet completely evaluated).

The concolic execution semantics is defined as the smallest relation that obeys the labeled transition rules shown in Figure 3. Given a trace \(\tau\), the initial symbolic state has the form \(\langle \tau; []; G_0; []; T \rangle\), where \(G_0\) is a goal with the same predicates as in the concrete execution, but with fresh variables as arguments, and \(T\)

---

¹ Nevertheless, an implementation of this technique may as well combine both calculi into a single algorithm to improve efficiency.

² The usefulness of keeping the clause stack will become clear in the next section. Informally speaking, it is needed to know which other clauses can be completely evaluated when a given clause—the one that is on top of the stack—is completely evaluated.
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<table>
<thead>
<tr>
<th>Rule</th>
<th>Description</th>
</tr>
</thead>
</table>
| (unfolding1) | $H \xrightarrow{\ell_1} B_n \ll P$ with $\text{mgu}(A_1, H_i) = \sigma, i = 1, \ldots, n$ and let $H \xrightarrow{\ell} B \ll P$ be a clause with $\ell \notin \{\ell_n\}$ and $\text{mgu}(A_1, H) = \sigma \neq \text{fail}$  
$$  \langle \text{\textit{u(\ell)}}(\ell_1)|\ell; |A_m, S; T\rangle \xrightarrow{s}\langle \text{\textit{c(\ell_1)}\sigma_1}\ell_2%; |(B_1, e(\ell_1), A_2, \ldots, A_m)\sigma_1; |(\ell_2); (B_2, e(\ell_2), A_2, \ldots, A_m)\sigma_2; \ldots |(\ell_n); (B_n, e(\ell_n), A_2, \ldots, A_m)\sigma_n]\rangle \ll S; T \rangle $$ |
| (unfolding2) | $H \xrightarrow{\ell_1} B_n \ll P$ with $\text{mgu}(A_1, H_i) = \sigma, i = 1, \ldots, n$ and there is no clause $H \xrightarrow{\ell} B \ll P$ with $\ell \notin \{\ell_n\}$ and $\text{mgu}(A_1, H) \neq \text{fail}$  
$$  \langle \text{\textit{u(\ell)}}(\ell_1)|\ell; |A_m, S; T\rangle \xrightarrow{s}\langle \text{\textit{c(\ell_1)}\sigma_1}\ell_2%; |(B_1, e(\ell_1), A_2, \ldots, A_m)\sigma_1; |(\ell_2); (B_2, e(\ell_2), A_2, \ldots, A_m)\sigma_2; \ldots |(\ell_n); (B_n, e(\ell_n), A_2, \ldots, A_m)\sigma_n]\rangle \ll S; T \rangle $$ |
| (exit) | $\langle \ell; (|L|); A; T; S; T \rangle \xrightarrow{\sigma} \langle \ell; A_2, \ldots, A_m; S; T \setminus \{\ell\} \rangle$ |
| (backtracking) | $\langle [b(\ell)]|\ell; |A_m; [\ell]; S; T\rangle \xrightarrow{\sigma} \langle \ell; |L; \sigma; S; T\rangle$ |
| (failure1) | there is $H \xrightarrow{\ell} B \ll P$ s.t. $\text{mgu}(A_1, H) = \sigma \neq \text{fail}$  
$$  \langle [f]|\ell; |A_m; [\ell]; T\rangle \xrightarrow{s}\langle \ell; \sigma; \text{fail}; [\ell]; T\rangle $$ |
| (failure2) | there is no $H \xrightarrow{\ell} B \ll P$ s.t. $\text{mgu}(A_1, H) \neq \text{fail}$  
$$  \langle [f]|\ell; |A_m; [\ell]; T\rangle \xrightarrow{s}\langle \ell; \sigma; \text{fail}; [\ell]; T\rangle $$ |

Fig. 3. Concolic execution semantics

is a set with the labels of all program clauses. The transition relation is labeled with the substitution computed in the step, and either a term $c(\ell, \theta)$ (a possible alternative for an unfolding step that concrete execution does not consider) or the special symbol $\perp$. Missing alternatives will be used to generated new input data that explore different execution paths.

Let us briefly explain the rules of concolic execution semantics:

- The unfolding1 rule follows the trace of the concrete execution and applies the same unfolding step. Here, a call of the form $e(\ell)$ is added to the end of the clause bodies to mark when the clauses are completely evaluated. This is required in our context since we only consider that a clause is covered when
all body atoms are successfully executed. This will be a useful information for test case generation, as we will see in the next section. Moreover, this rule considers that there is (at least) one more rule that unifies with the selected atom (we choose the first such clause). The label $\ell$ and the unifier $\sigma$ with this clause are also used to label the step. Moreover, we add $\ell_1$ to the stack of clause labels (the current environment).

- The unfolding rule proceeds similarly to the previous one but assumes that there are no further clauses whose head unifies with the selected atom.
- The exit rule applies when the selected atom has the form $e(\ell)$. In this case, we remove $\ell$ from the top of the environment stack, and also delete $\ell$ from the set of clause labels $T$ (i.e., clause $\ell$ has been completely evaluated).
- The backtracking rule proceeds similarly to the concrete semantics by applying a backtracking step according to the trace.
- Finally, the failure rules also proceed similarly to the concrete semantics but now we distinguish two cases, as in the unfolding rule. In particular, rule failure1 assumes that there is at least one clause whose head unifies with the selected atom, and rule failure2 considers that there is no such clause.

The rules unfolding and failure1 identify situations in which the symbolic state can follow an execution path that is not possible with the concrete goal. Therefore, they allow us to identify input data for the initial goal so that a different execution path can be followed.

Let us now show a simple computation with the concolic execution semantics.

We postpone to the next section the algorithm for test case generation.

Example 2. Consider the following simple program:

\[
\begin{align*}
(1) & \ p(X) :- q(X), r(X) . \\
(2) & \ q(X) :- s(X) . \\
(3) & \ s(a) . \\
(4) & \ s(b) . \\
(5) & \ r(b) . 
\end{align*}
\]

where we again consider natural numbers as clause labels. The concrete execution for the initial goal $p(a)$ is as follows:

\[
\langle p(a); [] \rangle \xrightarrow{u(1)} \langle q(a), r(a); [] \rangle \xrightarrow{u(2)} \langle s(a), r(a); [] \rangle \xrightarrow{u(3)} \langle r(a); [] \rangle \xrightarrow{\ell} \langle \text{fail}; [] \rangle
\]

Therefore, its associated trace is $\tau = [u(1), u(2), u(3), \ell]$. Now, for concolic execution, we consider the trace $\tau$ and the initial goal $p(X)$. The concolic execution is shown in Figure 4. As can be seen, the execution of clauses 1, 4 and 5 has not been completed. Moreover, we can observe that there was only one missing alternative when unfolding $s(X)$. In the next section, we show how this information can be used for test case generation.

\[3\] Observe that other, more relaxed, notions of clause covering are possible; e.g., consider that a clause is covered as soon as the clause is used in an unfolding step. Also, see [1] for a more declarative notion of test coverage.
Concolic Execution and Test Case Generation in Prolog

5 Test Case Generation

In this section, we present an algorithm for test case generation using concolic execution. In contrast to previous approaches for Prolog testing, our technique considers an underapproximation, i.e., only actual executions are considered (since there is no abstraction involved). Therefore, no false positives may occur. If a test case shows an error, this is an actual error in the considered program.

5.1 The Algorithm

In this section, we assume that the program contains a single predicate that starts the execution, which we denote with \textit{main}. This is not unusual for real applications. Moreover, we consider \textit{moded} logic programs. In particular, for simplicity, we assume that every predicate symbol \( p/n \) in the considered program has one associated mode,\(^4\) where \( \text{in}(p/n) = \{i_1, \ldots, i_m\} \) denotes the set of input parameters.

The algorithm for test case generation proceeds as follows:

1. First, a random goal of the form \( \text{main}(t_n) \) is produced, where at least the input arguments (according to \( \text{in}(\text{main}/n) \)) must be ground.
2. Now, we use the concrete semantics to execute the goal \( \text{main}(t_n) \), thus obtaining an associated trace \( \tau \). We assume that this execution terminates, which is reasonable since the input arguments are ground. In practice, one can use a timeout and report a warning when the execution takes more time.
3. Then, we use concolic execution to run an initial symbolic state of the form

\[
\langle [u(1), u(2), u(3), f]; [\text{}]; \ p(X); \ [\text{}]; [1, 2, 3, 4, 5] \rangle
\]

\[\vdash_{\text{id}}\]

\[
\langle [u(2), u(3), f]; [1]; \ q(X), r(X), e(1); \ [\text{}]; [1, 2, 3, 4, 5] \rangle
\]

\[\vdash_{\text{id}}\]

\[
\langle [u(3), f]; [2, 1]; \ s(X), e(2), r(X), e(1); \ [\text{}]; [1, 2, 3, 4, 5] \rangle
\]

\[
(\langle [f]; [3, 2, 1]; \ e(3), e(2), r(a), e(1); \ [\text{}]; [1, 2, 3, 4, 5] \rangle)
\]

\[\vdash_{\text{id}}\]

\[
\langle [f]; [2, 1]; \ e(2), r(a), e(1); \ [\text{}]; [1, 2, 4, 5] \rangle
\]

\[\vdash_{\text{id}}\]

\[
\langle [f]; [1]; \ r(a), e(1); \ [\text{}]; [1, 4, 5] \rangle
\]

\[\vdash_{\text{id}}\]

\[
\langle [\text{}]; [1]; \ \text{fail}; \ [\text{}]; [1, 4, 5] \rangle
\]

Fig. 4. Concolic execution for \([u(1), u(2), u(3), f]\) and \(p(X)\)

\(^4\) Extending our approach to multiple modes would not be difficult, but would introduce another source of nondeterminism when grounding an input goal.
the same steps). Let us consider that it has the following form:

\[
\langle \tau_0; L_0; G_0; S_0; T_0 \rangle \overset{c_1}{\sim} \langle \tau_1; L_1; G_1; S_1; T_1 \rangle \overset{c_2}{\sim} \langle \tau_2; L_2; G_2; S_2; T_2 \rangle \cdots \overset{c_m}{\sim} \langle \tau_m; L_m; G_m; S_m; T_m \rangle
\]

where \( \tau_0 = \tau, L_0 = \{\} \), \( G_0 = \text{main}(\overline{X_n}) \), \( S_0 = \{\} \), \( T_0 = T \), and \( G_m \) is either \text{true} or \text{fail}.

4. Now, we check the value of \( T_m \). If \( T_m = \{\} \), the algorithm terminates since all clauses have been completely executed. Otherwise, we identify the last state \( \langle \tau_i; L_i; G_i; S_i; T_i \rangle \) in the above concolic execution such that
   - the previous transition \( \overset{c_i}{\sim} \) is labeled with \( c_i = c(\ell, \theta_i) \), and
   - either \( \ell \in T_m \) or \( L_i \) contains (not necessarily in a top position) some labels from \( T_m \); the reason to also consider the labels from \( L_i \) is that considering an alternative clause may help to complete the execution of all the clauses in the current clause stack.

Therefore, we have a prefix of the complete concolic execution of the form:

\[
\langle \tau_0; L_0; G_0; S_0; T_0 \rangle \overset{c_1}{\sim} \langle \tau_1; L_1; G_1; \cdots \overset{c(\ell, \theta_i)}{\sim} \langle \tau_i; L_i; G_i; S_i; T_i \rangle
\]

and we are interested in the substitution \( \sigma_1 \sigma_2 \cdots \sigma_i \theta_i \), since it will allow us to explore a different execution path, possibly covering some more program clauses.

Hence, we have a second test case: \( G'_0 = \text{main}(\overline{X_n}) \sigma_1 \sigma_2 \cdots \sigma_i \theta_i \gamma \), where \( \gamma \) is a substitution that is only aimed at grounding the input parameters \( in(main/n) \) of \( \text{main} \) using arbitrary values (of the right type, preferably minimal ones).

5. Finally, we consider the initial state \( \langle G'_0, \{\} \rangle \) and obtain a new trace using the concrete execution semantics \( \tau' \), so that a new initial symbolic state is defined as follows:

\[
\langle \tau'; \{\}; \text{main}(\overline{X_n}); \{\}; T_m \rangle
\]

and the process starts again (i.e., we jump again to step 3). Observe that the initial state includes the set of clause labels \( T_m \) obtained in the last state of the previous concolic execution, in order to avoid producing new tests for clauses that are already covered by some previous test case.

Let us now illustrate the complete test case generation process with an example.

### 5.2 Test Case Generation in Practice

In this section, we illustrate the generation of test cases using a slight modification of the program in Example 1:

1. \( \text{main}(L,N,R) :- \text{length}(L,N), \text{rev}(L,[],R) \).
2. \( \text{main}(_-L,_,N,\text{error}) \).
3. \( \text{is_list}(\{\}) \).
4. \( \text{is_list}([H|T]) :- \text{is_list}(\{\}) \).
5. \( \text{is_list}(\{\}) \).

(1) main(L,N,R) :- length(L,N), rev(L,\[\],R).
(2) main(_L,_,N,\text{error}).
We start with a random initial goal, e.g., First iteration.

In some cases, there is no more goal instance of any previous goal, thus requiring some powerful abstraction operators. For instance, if there is a goal that includes an accumulating parameter, there are goals that are not instances since the search space is infinite and, moreover, due to the use of predicate \( \text{rev} \), hardly coincide. Also, using standard symbolic execution might be difficult too.

Consider the following prefix of the concolic execution:

\[
\langle [1, 2], u(8), b(2); [] \rangle; \text{main}(L, N, R); []; \{1, 2, 3, 4, 5, 6, 7, 8\}
\]

\[
\downarrow_{\text{id}} \quad \langle [u(8), b(2); [1]; \text{main}(L, N, R); []; \{1, \ldots, 8\})
\]

\[
\downarrow_{(b/\text{error})} \quad \langle [u(8), b(2); [1]; \text{main}(L, N, R); []; \{1, \ldots, 8\})
\]

\[
\downarrow_{c(r, (L/\{} N/0))} \quad \langle [u(8), b(2); [1]; \text{main}(L, N, R); []; \{1, \ldots, 8\})
\]

\[
\downarrow_{(b/\text{error})} \quad \langle [u(8), b(2); [1]; \text{main}(L, N, R); []; \{1, \ldots, 8\})
\]

Observe that, in this example, using a random generation of test cases would be useless since the length of the generated list and the second argument would hardly coincide. Also, using standard symbolic execution might be difficult too since the search space is infinite and, moreover, due to the use of predicate \( \text{rev} \) that includes an accumulating parameter, there are goals that are not instances of any previous goal, thus requiring some powerful abstraction operators.

Using concolic execution, though, we can easily generate appropriate test cases.

**First iteration.** We start with a random initial goal, e.g., \( \text{main}([a, b], s(0), R) \), where the input arguments \([1, 2]\) are assumed ground. The associated concrete execution is the following:

\[
\langle \text{main}(a, b), s(0), R); [] \rangle \quad \langle [u(1, 2), u(8)], \text{main}(L, N, R); []; \{1, \ldots, 8\})
\]

\[
\downarrow_{u(1, 2)} \quad \langle \text{main}(a, b), s(0), \text{rev}(a, b); [[2; \{R/\text{error}; \text{true}\}])
\]

\[
\downarrow_{u(8)} \quad \langle \text{main}(a, b), s(0), \text{rev}(a, b); [[2; \{R/\text{error}; \text{true}\}])
\]

and its associated trace is thus \( \tau = [u(1, 2), u(8), b(2)] \).

Now, we use concolic execution and produce the computation shown in Figure 5. Therefore, by executing \( \text{main}([a, b], s(0), R) \) only clause (2) is completely evaluated. According to the previous algorithm for test case generation, we now consider the following prefix of the concolic execution:

\[
\langle [u(1, 2), u(8), b(2); [1]; \text{main}(L, N, R); []; \{1, \ldots, 8\}) \quad \downarrow_{\text{id}} \quad \langle [c(r, (L/\{} N/0)) \rangle \quad \ldots
\]

and the associated substitution \([L/\{} N/0\}).
Fig. 6. Concolic execution for \( \langle u(1, 2), u(7), u(3) \rangle \); \([\ ]; \text{main}(L, N, R); [\ ]; \{1, 3, \ldots, 8\}\)

**Second iteration.** In this case, no further grounding of the goal is required, and we start a new concrete execution for the goal \text{main}(\[\ ], 0, R):

\[
\langle \text{main}(\[\ ], 0, R); [\ ] \mid u(1,2) \rangle \\
\langle \text{length}(\[\ ], 0), \text{rev}(\[\ ], [], R); [2; \{R/error\}; \text{true}] \rangle
\]

\[
u(2) \rightarrow \langle \text{rev}(\[\ ], [], R); [2; \{R/error\}; \text{true}] \rangle
\]

\[
u(3) \rightarrow \langle \text{true}; [2; \{R/error\}; \text{true}] \rangle
\]

and obtain the following trace \( \tau' = \langle u(1, 2), u(7), u(3) \rangle \). Now, we perform the concolic execution shown in Figure 6. Therefore, according to the algorithm, we consider the following prefix of the concolic execution:

\[
\langle [u(1, 2), \ldots; [\ ]; \text{main}(L, N, R); [\ ]; \{1, 3, \ldots, 8\}) \rightarrow \langle \ldots \rangle \stackrel{c(L/X, N; s(N'))}{\rightarrow} \langle \ldots \rangle
\]

with the associated substitution \{L/X=L', N/s(N')\}.

**Third iteration.** Now, we consider the goal \text{main}(X|L', s(N'), R). Since the first two arguments must be ground, as mentioned before, we apply a minimal grounding substitution and get, e.g., \text{main}(a, s(0), R). The concrete execution, which is shown in Figure 7, computes the following trace:

\[
\tau'' = \langle u(1, 2), u(8), u(7), u(4), u(5), u(3) \rangle
\]

Then, we use concolic execution again as shown in Figure 8. Therefore, according to the algorithm, we consider the following prefix of the concolic execution:

\[
\langle [u(1, 2), u(8), \ldots; [\ ]; \text{main}(L, N, R); [\ ]; \{4, 5, 6, 8\}) \rightarrow \langle \ldots \rangle \stackrel{c(L/X, N; s(N'))}{\rightarrow} \langle \ldots \rangle \stackrel{c(L/X, N; s(N'))}{\rightarrow} \langle \ldots \rangle
\]

with the associated substitution \( \sigma_1 \theta_2 = \{L/X=Y|L''', N/s(N'')\} \).

**Fourth (and last) iteration** As in the previous case, the instantiated goal, \text{main}(X, Y|L'''', s(s(N'''), R), is not ground enough according to its input mode.
and, thus, we apply a minimal grounding substitution. In this case, we get the initial goal main([a, b], s(s(0)), R). Here, the concrete execution is basically the same shown in Figure 2, except for the last (backtracking) step. Therefore, the associated trace is

$$\tau'' = [u(1, 2), u(8), u(8), u(7), u(4), u(5), u(4), u(6), u(5), u(3)]$$

Now, concolic execution from the initial state

$$([u(1, 2), u(8), u(8), u(7), u(4), u(5), u(4), u(6), u(5), u(3)]; []; \text{main}([\_], N, R); []; \{6\})$$

proceeds similarly to the derivation shown in Figure 8, but now clause (6) is also completely evaluated, which means that the algorithm terminates successfully.

To summarize, concolic testing generated four test cases:

- main([a, b], s(0), R)
- main([\_], 0, R)
- main([a, s(0)], R)
- main([a, b], s(s(0)), R)

which suffice to cover the complete evaluation of all program clauses.

In general, when the test case generation algorithm terminates, concolic testing is sound (i.e., there are no false positives since only concrete executions are considered) and complete (in the sense that all clauses are completely evaluated when using the computed test cases, i.e., we get a 100% coverage). When the process is stopped (e.g., because it does not terminate or takes too much time), our test case generation is only sound. Note that this contrasts with other approaches to test case generation in Prolog (and CLP), e.g., [7, 14], where full coverage is not considered.

Formally proving these results is not difficult. The challenge, however, is experimentally verifying the effectiveness and scalability of our approach with real Prolog programs. For this purpose, though, we first need to extend concolic execution to deal with negation, built-in’s, extra-logical features, etc. This extension
The implementation has been undertaken with promising results. An approach to Prolog testing and debugging, which we plan to extend to deal with full Prolog, has been shown quite successful in other programming paradigms, especially when dealing with large applications. Therefore, it might have a great potential for Prolog testing, too.

We have introduced a novel approach to Prolog testing and debugging. The so-called concolic execution that mixes concrete and symbolic execution has been shown quite successful in other programming paradigms, especially when dealing with large applications. Therefore, it might have a great potential for Prolog testing, too.

In this paper, we have considered a limited scenario: pure Prolog without negation. Nevertheless, the main distinctive features of the Prolog programming language—i.e., unification, non-determinism and backtracking—are present here, so it was not trivial at all to adapt standard concolic execution [5, 15] to Prolog. Now, we plan to extend this approach to deal with full Prolog. Dealing with arithmetic built-ins, for instance, can be done by producing a path constraint, analogously to symbolic execution in imperative languages. An implementation has been undertaken with promising results.

6 Concluding Remarks and Future Work

is the subject of ongoing work. For this purpose, we will consider the linear operational semantics of [16], and its symbolic version [4], as a promising starting point.
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References

Liveness properties in CafeOBJ – a case study for meta-level specifications*

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Abstract. We provide an innovative development of algebraic specifications and proof scores in CafeOBJ by extending a base specification to the meta-level that includes infinite transition sequences. The infinite transition sequences are modeled using behavioral specifications with hidden sort, and make it possible to prove safety and liveness properties in a uniform way.

As an example of the development, we present a specification of Dijkstra’s binary semaphore, a protocol to guarantee exclusive access to a resource. For this protocol we will give three different properties, one being the mutual exclusion (or safety) property, and two more regarding different forms of liveness, which we call progress property and entrance property. These three properties are verified in a computationally uniform way (by term rewriting) based on the new development.

Besides being a case study of modeling meta-properties in CafeOBJ, we provide an initial characterization of strength of various properties. Furthermore, this method can serve as a blue-print for other specifications, in particular those based on Abstract State Systems (Ass).

Keywords: algebraic specification, liveness, CafeOBJ, verification

1 Introduction

QLOCK, an abstract version of Dijkstra’s binary semaphore, is a protocol to guarantee exclusive access to a resource. Besides the initial specification and verification in CafeOBJ (see for example [2]), it saw implementations in Coq [7] and MAUDE [9]. Most of these specifications only consider safety properties, in the current case the mutual exclusion property, that no two agents will have access to the resource at the same time. However, liveness properties are normally left open. These properties ensure that ‘there is progress’. In our particular case, they ensure that agents do not block out other agents from acquiring access to the resource.

We are using CafeOBJ as specification and verification language. CafeOBJ is a many- and order-sorted algebraic specification language from the OBJ family,

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related to languages like Casl and Maude. CAFEOBJ allows us to have both the specification and the verification in the same language. It is based on powerful logical foundations (order-sorted algebra, hidden algebra, and rewriting logic) with an executable semantics \([3, 5, 6]\).

The particular interest of the current development is two-fold: Firstly, it extends base specifications in order-sorted and rewriting logics to a meta-level, which requires behavioral logic, thus using the three logics together to achieve the proofs. Secondly, we use a search predicate and covering state patterns that allow us to prove the validity of a property over all possible one-step transitions, by which safety and liveness properties in the base and meta-level can be proven.

*Layout of the article* In Section 2 we introduce the QLOCK protocol and various properties for verification, give a short introduction to the CAFEOBJ language, and provide the base specification onto which the current work is building. This section also discusses briefly the proof method by induction and exhaustive search.

In Section 3 we extend the base specification to include infinite transition sequences. Here we also discuss the methodology of meta-modeling.

In Section 4 we provide (parts) of the proof score verifying the three properties.

In the final section we provide a discussion of the approach with respect to applicability to different problems, and conclude with future research directions.

## 2 The QLOCK protocol

The QLOCK protocol regulates access of an arbitrary number of agents to a resource by providing a queue (first-in-first-out list). Agents start in the remainder section, henceforth indicated by rs. The mode of operation is regulated by the following set of rules:

- If an agent wants to use the resource, it puts a unique identifier into the queue, and by this it transitions into the waiting section (ws).
- In the waiting section, an agent checks the top of the queue. If it is the agent’s unique identifier, the agent transitions into the critical section (cs), during which the agent can use the resource.
- After having finished with the resource usage, the agent removes the head of the queue and transitions back into the remainder section (rs).

See Figure 1 for a schematic flow diagram.

### 2.1 Verification properties

The basic safety property of QLOCK is the *mutual exclusion property* (mp):

**Property 1 (mutual exclusion property).** At any time, at most one agent is in the critical section.
While this is the most important property for safety concerns, it does not guarantee that an agent wanting to use the resource ever gets the chance to use it (for example, in case of denial-of-service attack to a server). To guarantee this, we define two liveness properties: The first concerns the transition from ws to cs, and is called the progress property (pp). This property has already been discussed in [8] as lockout freedom property.

Property 2 (progress property). An agent that has entered into the waiting section (ws), i.e., has put his unique identifier into the queue, will eventually transition into the critical section (cs), i.e., progress to the top of the queue and gain access to the resource.

The last one concerns the transition from the remainder section rs to the waiting section ws, called entrance property (ep):

Property 3 (entrance property). An agent will eventually transition into the queue, i.e., from the remainder section to the waiting section.

Although it might sound counter-intuitive that the entrance property should hold for each agent at all times, we believe that there are good reasons to consider this property: This is motivated by the fact that given any finite run, i.e., finite transition sequence, of the QLOCK protocol, we can always extend it to an infinite and fair transition sequence (see later sections for details). Thus, what we are actually proving is that for each agent, either the entrance property holds, or the execution terminates before the agent had a chance. In circumstances of long-running services (like most client-server interaction where a server is practically never stopped), this ensures that agents, or clients, will – given long enough execution time – eventually be served.

As we will see later on, to prove this property we need additional assumptions, in particular fairness, see Section 3.2, which makes it conceptually different from the first two properties. To continue with the analogy set forth above, the schedulers used in most operating systems or network hubs try to create a fair execution sequence by using round-robin or similar techniques [14,15].
2.2 Short introduction to CafeOBJ

Although we cannot give a full introduction to the CafeOBJ language, to aid readers unfamiliar with it we give a short introduction. Users acquainted with Maude can safely skip this section, as syntax of the two languages is very similar.

CafeOBJ is an algebraic specification language, thus the unit of specification is an algebra, in particular an order-sorted algebra. To specify an algebra, the following information have to be given:

**signature** Similar to normal algebras, a signature consists of operators and their arities. In the multi-sorted setting we are working in, this means that sorts have to be defined, and for each operator (or function) the number and sorts of the arguments and result have to be specified.

**axioms** To create more than the free algebra, we need axioms (or equations) between terms.

We will demonstrate these concepts on a simple definition of natural numbers:

```plaintext
mod! SIMPLE-NAT {
  signature {
    [ Zero NzNat < Nat ]
    op 0 : -> Zero
    op s : Nat -> NzNat
    op _+_ : Nat Nat -> Nat
  }
  axioms {
    vars N N' : Nat
    eq 0 + N = N .
    eq s(N) + N' = s(N + N') .
  }
}
```

Line 1 begins the specification of the algebra called SIMPLE-NAT with initial semantics (indicated by the ! after mod). The body of the specification consists of the above mentioned two blocks. Lines 2-7 define the signature, lines 8-12 the axioms. In line 3 the sorts and their order is introduced. In this case there are three sorts, Zero, NzNat, and Nat. The order relation gives us that the former two are a sub-sort of Nat, but does not specify a relation between themselves. Lines 4-6 give three operators, one being the constant 0 of sort Zero, one being the successor function s, and one being addition, written in infix notation. Here the _ are the places of the arguments.

The second block defines the equations by first declaring two variables of sort Nat. Axioms are introduced with the eq keyword, and the left- and right-hand side of the equation are separated by =. Thus, the two axioms provided here provide the default inductive definition of addition by the successor function.

In the following, the signature and axioms block declaration will be dropped, as they are not necessary.
2.3 Base specification

We are building upon a previously obtained specification and verification of QLOCK [2]. Although we are providing the complete source of the specification part, we cannot, due to space limits, include the full verification part. The reader is referred to the full code at [11].

The basic idea of the following specification is to use the natural abstraction of QLOCK and its transitions as an Abstract State System (Ass). The usual steps in providing an algebraic specification in this setting are:

– Define proper abstraction as modules/algebras of the players (e.g., agent, queue) in the protocol.
– Model the complete state as a module/algebra.
– Use transitions between terms of the state algebra to describe transitions in the protocol.
– Specify (and later verify) properties on states (and transitions) to ensure safety and liveness.

Let us start with the most basic item, namely modeling the sections an agent can be in. For literals, i.e., elements of the sort LabelLt, which is a sub-sort of the sort Label, we define identity via the syntactical identity. The second module specification defines an agent, or more specifically, the algebra of agent identifiers AID, without any further axioms, which implies that identifiers are considered different if they are syntactically different.

```
mod! LABEL { [ LabelLt < Label ]
  vars L1 L2 : LabelLt .
  ops rs ws cs : -> LabelLt {constr} .
  eq (L1 = L2) = (L1 == L2) . }
mod* AID { [ Aid ] }
```

In the next step we model a queue, a first-in-first-out storage. Note in the following code, that CAFEOBJ allows for parametrized modules. In the present case the parameter X has no further requirements, which is expressed by the fact that it only needs to belong to the trivial algebra. Another important point to note is that we are using associative constructors, which allows us to freely use any way of parenthesizing. Similarly, we introduce a module for parametrized sets, where we use associative and commutative constructors.

```
mod! QUEUE (X :: TRIV) { [ Elt.X < Qu ]
  vars Q Q1 Q2 : Qu . vars E E1 E2 : Elt .
  op empQ : -> Qu {constr} .
  op (_&_) : Qu Qu -> Qu {constr assoc id: empQ} .
  eq (empQ = (E & Q)) = false .
  eq ((E1 & Q1) = (E2 & Q2)) = ((E1 = E2) and (Q1 = Q2)) . }
mod! SET(X :: TRIV) { [ Elt.X < Set ]
  vars E : Elt .
  op empty : -> Set {constr} .
  op (_,_) : Set Set -> Set {constr assoc comm id: empty} .
  eq E E = E . }
```
Concerning agents, we model them as terms of an algebra of *agent observers* which associates agent identifiers with labels, expressing the fact that the agent is in the current state. More formally, the meaning of the term \( \text{lb}[A]:S \) is that the agent \( A \) is in section \( S \):

\[
\text{mod! AOB \{protecting\( (LABEL) \) protecting\( (AID) \) \[ Aob \]}
\quad \text{op} \ (\text{lb[\_]:\_}) : \text{Aid Label} \rightarrow \text{Aob \{constr\} .}
\]

In the final step we instantiate the parametrized queue with agent ids, and define the state algebra as a pair of one queue and an arbitrary set of agent observers. Note that the pairing is done by the syntax \( 1 \$ r \), CAFEOBJ allows nearly arbitrary syntax:

\[
\text{mod! AID-QUEUE \{ protecting\( (\text{QUEUE}(\text{AID}\{\text{sort Elt} \rightarrow \text{Aid}\}) ) \) \}}
\]
\[
\text{mod! STATE\{ protecting\( (\text{AID-QUEUE}) \}
\quad \text{protecting\( (\text{SET}(\text{AOB}\{\text{sort Elt} \rightarrow \text{Aob}\})\{\text{sort Set} \rightarrow \text{Aobs}\})\)
\[ \text{State} \] \text{op} \ \_\_ : \text{Qu Aobs} \rightarrow \text{State \{constr\} .}
\]

With this we have given a complete definition of the state algebra, but the dynamic aspect of the protocol has been left out till now. We are now providing transition rules over states to express this dynamic aspect. In the following code segments, the two states of the transition are aligned, and changing parts are indicated with a bold font. The three transitions are \text{WaitTrans}, where an agent transitions from \( rs \) to \( ws \), \text{TryTrans}, where an agent tries to enter \( cs \), and \text{ExitTrans}, where an agent leaves the critical state:

\[
\text{mod! WaitTrans \{ protecting\( (\text{STATE}) \).} \\
\quad \text{var Q : Qu . var A : Aid . var AS : Aobs .} \\
\quad \text{trans[wt]}: \ (Q \quad \$ \ ((\text{lb}[A]: \text{rs}) \ AS)) \\
\quad \quad \Rightarrow \ ((Q \ & \ A) \quad \$ \ ((\text{lb}[A]: \text{ws}) \ AS)) . \}
\]
\[
\text{mod! TryTrans \{ protecting\( (\text{STATE}) \).} \\
\quad \text{var Q : Qu . var A : Aid . var AS : Aobs .} \\
\quad \text{trans[ty]}: \ ((A \ & \ Q) \quad \$ \ ((\text{lb}[A]: \text{ws}) \ AS)) \\
\quad \quad \Rightarrow \ ((A \ & \ Q) \quad \$ \ ((\text{lb}[A]: \text{cs}) \ AS)) . \}
\]
\[
\text{mod! ExitTrans \{ protecting\( (\text{STATE}) \).} \\
\quad \text{var Q : Qu . vars A1 A2 : Aid . var AS : Aobs .} \\
\quad \text{trans[ex]}: \ ((A1 \ & \ Q) \quad \$ \ ((\text{lb}[A2]: \text{cs}) \ AS)) \\
\quad \quad \Rightarrow \ (Q \quad \$ \ ((\text{lb}[A2]: \text{rs}) \ AS)) . \}
\]

Based on the above specification, it is possible to provide a *proof score*, i.e., a program in CAFEOBJ, that verifies the mutual exclusion property \( \text{mp} \). As usual with proofs by induction over the reachable states (see below), the target property by itself does not suffice to work as inductive property, making the introduction of further properties on states necessary. Example properties that have to be used are uniqueness properties (e.g., the same agent identifier cannot appear several times in the queue) or initial state properties (e.g., the queue is empty at the beginning). Obtaining an inductive property (set of properties) is one of the challenging aspects of verifications, and requires an iterative and interactive approach. Readers interested in the details are referred to the code at [11].
We conclude this section with a short discussion on the verification methodology applied here. Details concerning this verification technique and a more generalized methodology will be presented at [12] and a forthcoming article.

2.4 Verification by induction and exhaustive search

Verification of properties of an Ass is often done by induction on reachable states, more specifically by induction over the length of transition sequences from initial states to reachable states. That is, we show that a certain property (invprop) holds in the set of initial states, characterized by init. Furthermore, as we proceed through transitions (state changes), the invprop is preserved.

But to show liveness properties, considering only invariant properties on states is not enough. We thus use an extended method that does inductive proofs on the reachable state space, and in parallel proves properties (transprop) on all transitions between reachable states. To be a bit more specific, assume that $S \Rightarrow S'$ is a transition from one state (state term, state pattern) $S$ to a state $S'$. We show that if $\text{invprop}(S)$ holds, then also $\text{invprop}(S')$ (the induction on reachable states), but also that for this transition $\text{transprop}(S, S')$ holds.

Both of these are done with CAFEOBJ's built-in search predicate (see Section 4.2), which exhaustively searches and tests all possible transitions from a given state (pattern). The concepts introduced here are an extension and generalization of transition invariants [10], details in a forthcoming publication.

In the CAFEOBJ setting, which means rewrite-based, we have to ensure that both of the following implications reduce to True:

\[
\text{init}(S) \rightarrow \text{invprop}(S) \\
\text{invprop}(S) \rightarrow \text{invprop}(S') \quad \text{where } S \rightarrow S' \text{ is a transition}
\]

where $S$ and $S'$ are states (state terms) describing the pre- and post-transition states, respectively. This has to be checked for all possible transitions available in the specification.

If this can be achieved, we can be sure that in all reachable states, i.e., those that can actually occur when starting from an initial state, the required property invprop holds.

3 Extended specification

The starting point of the following discussion is the question of how to verify liveness properties. Initial work by the first author led to a specification which kept track of the waiting time in the queue. Combined with the assumption that there are only finitely (but arbitrary) many agents, we could give a proof score not only for the mutual exclusion property mp, but also for the progress property pp. This work was extended by the third author to the currently used base specification.

To verify the last property, ep, operational considerations alone do not suffice. On the level of observers, we cannot guarantee that an agent will ever enter
the queue, since we have no control over which transitions are executed by the
system. To discuss (verify) this property, we have to assume a certain meta-level
property, in this case the fairness of the transition sequence. A similar approach
has been taken in [4] for the Alternating Bit Protocol, where fair *event mark
streams* are considered.

3.1 Fairness

The concept of fairness we are employing here is based on the mathematically
most general concept:

**Definition 1 (Fairness).** A sequence of transitions $S$ is called fair, if every
finite sequence of transitions appears as sub-sequence of $S$.

A necessary consequence of this definition is that every fair sequence is infinite.

**Relation to other concepts of fairness** The methodology of using Ass in
CAFEOBJ has been strongly influenced by UNITY [1], which builds upon a
semantics similar to Ass of sequences of states and temporal operators. It provides
an operator *ensures*, which can be used to model fairness via a measure function.

Another approach to the concept of fairness is taken by LTL logic [13], where
two types of fairness, strong and weak, are considered, referring to *enabled* and
*applied* state of transitions.

\[
\text{weak } \Diamond \square \text{enabled}(t) \rightarrow \square \Diamond \text{applied}(t) \\
\text{strong } \square \Diamond \text{enabled}(t) \rightarrow \square \Diamond \text{applied}(t)
\]

where *enabled* and *applied* are properties on transition instances. In the particular
case we are considering, *enabled* is always true, as we can execute every instance
of a transition at any time, due to the fact that the wait-state transition can
be applied even if the agent is not at the top of the queue. Fairness in this case
means that, at some point every transition will be applied.

Both concepts can be represented in suitable way by the definition of fairness,
or in other words, the definition used in this setting (every finite sequence is
sub-sequence) subsumes these two concepts.

3.2 Transition sequence

As mentioned above, modeling fairness requires recurring to a meta-assumption,
namely that the sequence of transitions is fair, i.e., every instance of a transition
appears infinitely often in the sequence. In our case we wanted to have a formal-
ization of this meta-assumption that can be expressed with the rewriting logic of
CAFEOBJ.

The approach we took models transition sequences using behavioral specifica-
tion with hidden algebra [4], often used to express infinite entities. Note that we
are modeling the transition sequence by an infinite stream of agent identifiers,
since the agent uniquely defines the instance of transition to be used, depending
on the current state of the agent. This is a fortunate consequence of the modeling
scheme at the base level, where, if we pick an agent, we can always apply the
transition that is uniquely defined by that agent. Translated into the language
of the above mentioned LTL logic it means that all transitions are permanently
enabled.

\[ mod* \text{TRANSSEQ} \{ \text{protecting(AID)} \]
\[ *[\text{TransSeq}]* \]
\[ op (_&_) : \text{Aid TransSeq} \rightarrow \text{TransSeq} . \]

The transition sequence is then used to model a meta-state, i.e., the combi-
nation of the original state of the system as specified in the base case, together
with the list of upcoming transitions:

\[ mod! \text{METASTATE} \{ \text{protecting(STATE + ... ) [MetaState]} \]
\[ op _-_ : \text{State TransSeq} \rightarrow \text{MetaState} \{\text{constr}\} . ... \]

The dots at the end of the definition refer to a long list of functions on
meta-states, we will return to this later on.

In the same way, transitions from the base case are lifted to the meta-level.
Here we have to ensure that the semantics of transition sequences and the
transition in the original (non-meta level) system do not digress. That means
first and foremost, that only the agent at the top of the transition sequence can
be involved in a transition.

Let us consider the first transition WaitTrans:

\[ mod! \text{MWT} \{ \text{protecting(METASTATE)} \var Q : \text{Queue} . \]
\[ \var A : \text{Aid} . \var AS : \text{Aobs} . \var T : \text{TransSeq} . \]
\[ \text{trans [meta-wt]}: \]
\[ ( \text{Q } \&( \text{lb[ A ]}: \text{rs} AS)) ^ (A & T) ) \]
\[ => ( \text{Q & A } \&( \text{lb[ A ]}: \text{ws} AS) ^ T) . \]

Due to the structural definition of the meta-transition, we see that it can only
be applied under the following conditions that

- the agent is at the top of the transition sequence, and
- the agent is in remainder section rs.

The next transition is TryTrans, where an agent checks whether it is at the
top of the queue, and if yes, enters into the critical section. Adding the meta-
level requires only that the agent is also at the head of the transition sequence.
This transition uses the built-in operator \text{if_then_else_fi}, because we want to
destructively use up the top element of the transition sequence to ensure that no
empty transitions appear.

\[ mod! \text{MTY} \{ \text{pr(METASTATE)} \var Q : \text{Queue} . \]
\[ \text{vars A B : Aid} . \var AS : \text{Aobs} . \var T : \text{TransSeq} . \]
\[ \text{trans [meta-ty]}: \]
\[ (((B & Q) \&( \text{lb[ A ]}: \text{ws} AS)) ^ (A & T)) \]
=> if (A = B) then
    (((A & Q) $ ((lb[ A ]: cs) AS)) ~ T)
else
    (((B & Q) $ ((lb[ A ]: ws) AS)) ~ T)
fi.
}

The final transition is ExitTrans, where an agent returns into the remainder section:

mod! MEX {pr(METASTATE) var Q : Queue .
  var A : Aid . var AS : Aobs . var T : TransSeq .
  trans[meta-ex]:
    (((A & Q) $ ((lb[ A ]: cs) AS)) ~ (A & T))
  => (((Q $ ((lb[ A ]: rs) AS)) ~ T).
}

Combining all these algebras provides the specification of the meta system:

mod! METAQLOCKsys{ pr(MWT + MTY + MEX) }

As mentioned above, to express the fairness condition, we recur to an equivalent definition, namely that every finite sequence of agent identifiers can be found as a sub-sequence of the transition sequence, rephrased here in an indirect way in the sense that it cannot happen that we cannot find a (finite) sequence of agent identifiers in a transition sequence:

eq ( find ( Q, T ) = empQ ) = false .

### 3.3 Waiting times

The axiom shown above uses the function find, which has been defined in the algebra METASTATE. We mentioned that several additional functions are defined, too. These functions are necessary to compute the waiting time for each agent.

Waiting time here refers to the number of meta-transitions until a particular agent is actually changing its section. Here complications arise due to the trial transition from ws to cs (by means of the if_then_else_fi usage), where an agent might find itself still in ws due to not being at the head of the queue. This has to be considered while computing waiting times for each agent.

Closer inspection of the transition system provides the following values for waiting times:

For agents in rs and cs The waiting time is determined by the next occurrence of its agent id in the transition sequence, since there are no further requirements. In CAPEOBJ notation, the length of the following sequence:

```
find( A, T )
```
Liveness properties of in CafeOBJ

Fig. 2. Computing of waiting time from ws to cs

For agents in ws Here we have to ensure that the agent advances to the top of the queue. Thus, each agent higher up in the queue has to appear two times in the transition sequence, once for entering the critical section, and once for leaving it. Let \( Q_A \) be the part of the queue that is above (higher up) the agent id \( a \). Then the waiting time for \( a \) would be determined by doubling \( Q_A \), then searching the transition sequence first for the doubled \( Q_A \), and finally searching for the next appearance of \( a \). Consider for example the state represented in Figure 2. Assume that initially the queue contains the three agent identifiers \( b \), \( c \), and \( a \) (in this order), and all of them are initially in rs. To see when the \( a \) at the bottom of the queue can transition into cs, we first have to bring \( b \) into cs, which happens at position 2. After that there is a series of trial transitions without success (\( a \), \( c \), \( a \)) until another \( b \) appears, which makes the agent \( b \) transition back into rs. At this point the queue contains two elements, \( c \) and \( a \). The same repeats for \( c \), until finally \( a \) can enter into cs. Summing up, what has to be searched within the transition sequence is the sub-sequence \( b b c c a \) which amounts to

\[
\text{find} \left( \text{double} \left( \text{droplast} \left( \text{find} ( A, Q ) \right) \right) \& A, T \right)
\]

(The actual code is slightly different, but semantically the same.)

Here the doubling of the \( Q_A \) is achieved by first finding the sub-sequence within the queue up to \( A \) (which includes \( A \)), and then dropping the last element, which is \( A \).

The functions mentioned above are exactly those needed to compute this waittime function. Due to the necessity of error handling, the definition becomes a bit lengthy.

4 Verification of properties

Our aim is to verify the progress property and the entrance property. These properties are now expressed in the following ways:

– At any transition, the waiting time of agents not changing section decreases.
– If the waiting time of an agent reaches 0, then a section change happens.

Combining these two properties, and assuming fairness, we can show both, that every agent will eventually enter into the queue, and every agent in the queue will eventually gain access to the resource, i.e., enter into the critical section.

In the following let us assume that the following variable definitions are in effect:
Then the CAFEOBJ implementation of the first property is as follows:

pred wtd-allaid : Aobs MetaState MetaState .
eq[::and wtd-allaid]:
  wtd-allaid( ( ( lb[A]: C ) AS ) , S, SS ) =
   ( ( sec( ( lb[A]: C ) ,S) == sec( ( lb[A]: C ) , SS) ) implies
     ( waittime( A, S ) > waittime ( A, SS )) ) .

And the one of the second property:

pred wtzerochange-allaid : Aobs MetaState MetaState .
vars S SS : MetaState . var Q : Queue . var AS : Aobs .
var A : Aid . var C : Label . var QQ : TransSeq .
eq[::and wtzerochange-allaid]:
  wtzerochange-allaid( ( ( lb[ A ]): C ) AS ) , S , SS ) =
   ( ( waittime( A, S ) == 0 ) implies
     ( sec( ( lb[ A ]): C ) , S ) /= sec( ( lb[ A ]): C ) , SS ))) .

Here we have to note that the sec operator computes the actual section SS and not the one given by C.

These two properties alone do not function as inductive invariant, so several more have to be included. In addition, we are lifting also the properties used in the original specification to the meta level by making the new operators simply operate on the projections. Again, the interested reader is referred to [11] for the full source.

4.1 Proof score with patterns

The method described in Section 2.4 is used here in combination with a covering set of patterns. We mention only the definition of cover set here, but details on this methodology will be presented at [12] and a forthcoming article:

**Definition 2 (cover set).** Assume a set $S \subseteq \text{State}$ of states (ground terms of sort State) is given. A finite set of state patterns $C = \{C_1, \ldots, C_n\}$ is called cover set for $S$ if for every $s \in S$ there is a substitution $\delta$ from the variables $X$ occurring in $C$ to the set of all ground terms, and a state pattern $C \in C$ such that $\delta(C) = s$.

Practically this means, that we give a set of state terms that need to cover all possible ground instances of state terms. For the base case, a set of 13 state patterns has been given. We list only the cases for rs, the cases for the other sections are parallel.

```
eq s1 = (q $ empty) .
eq s2 = (empQ $ ((lb[b1]: rs) as)) .
eq s8 = ((b1 & q) $ ((lb[b1]: rs) as)) .
eq s11 = ((b1 & q) $ ((lb[b2]: rs) as)) .
```
For the meta-level we combine these patterns with patterns for the transition sequence, where once \( b_1 \) is at the head of the transition sequence, and once another identifier \( b_2 \), amounting to 26 different meta state patterns:

\[
\begin{align*}
eq n_1 &= (s_1 \neg (b_1 \& t)) \\
eq n_2 &= (s_2 \neg (b_1 \& t)) \\
eq l_1 &= (s_1 \neg (b_2 \& t)) \\
eq l_2 &= (s_2 \neg (b_2 \& t)) \\
&\ldots
\end{align*}
\]

We conclude this section with a discussion of the search predicate, actually family of search predicates, in \textsc{CafeOBJ}.

### 4.2 The \textsc{CafeOBJ} search predicate

During a proof over reachable states by induction on the transitions, we need a method that provides \textit{all} possible successors of a certain state. The \textsc{CafeOBJ} search predicate we use is \( S =(*,1)=>+ S' \) suchThat \( prop(S, S') \), where \( S \) and \( S' \) are states, and \( prop(S, S') \) is a Boolean value. This is a tricky predicate and full discussion is available in an upcoming reference manual for \textsc{CafeOBJ}, but in this case it does the following:

- It searches all successor states \( S' \) that are reachable in exactly one step from the left side state \( S \) (here 1 stands for maximal one step, and + for at least one step). Call the set of successors \( Succ(S) \).
- It checks all those states determined by the first step whether the property given in \textsc{Bool} holds. If there is at least one successor state where it holds, the whole predicate returns true, i.e., what is returned is \( \exists S' \in Succ(S) : prop(S, S') \).

This can be used with a double negation to implement an exhaustive search in all successor states by using the equality:

\[
\forall S' \in Succ(S) : prop(S, S') \iff \neg \exists S' \in Succ(S) : \neg prop(S, S')
\]

This leads to the following, admittedly not very beautiful, definition of the inductive invariant condition, where we use \( SS \) for the \( S' \) in the above equality:

\[
\begin{align*}
\text{pred inv-condition : MetaState MetaState .} \\
\text{eq inv-condition(S:MetaState, SS:MetaState) =} \\
\text{ ( not ( S =(*,1)=>+ SS suchThat} \}
\text{ ( not ( inv-prop(S, SS) == true))))} \\
\end{align*}
\]

Here \text{inv-prop} is the set of inductive invariant properties we have mentioned above.

Note that the operator used here has access not only to the original or the successor state, but to both states. This peculiar feature allows us to prove properties like decrease of waiting time, which is impossible if there is no access to both values in the same predicate. As pointed out above, \textsc{CafeOBJ} actually includes a whole set of search predicates which allows searching for arbitrary, but given depth, but verifications making use of these predicates are still to come.

The final step in the proof score is to verify that both the initial condition and the inductive invariant condition do actually hold on all the state patterns by reducing the expressions to true:
red init-condition(n1) . red init-condition(n2) . ...
red inv-condition(n1, SS: MetaState) .
red inv-condition(n2, SS: MetaState) . ....

This concludes the discussion of the specification methodology and proof score code.

5 Discussion and conclusion

Using the method laid out above, we have formally verified the three properties given in the beginning, mutual exclusion property (only at most one agent at a time is in the critical section), progress property (an agent in the queue will eventually gain access to the resource), and entrance property (every agent will eventually enter into the queue). While the original base specification and proof score verified the first two properties, it also required the assumption that the number of agents is finite. In our case, this assumption is superseded by the assumption of fairness of the transition sequence, which is by itself necessary to verify the third property.

This methodology also opens up further options in the specification. By requiring acquisition of the resource within a reasonable time, and providing requirements on the transition sequence that the reasonable-time condition is fulfilled, we believe it is possible to specify and verify time-critical systems.

We have to note that what we call here progress property has already been shown in different settings [8]. The key contribution is the extension to the entrance property, meaning that an agent always gets a chance to enter the queue. In addition, we extended the proof of the progress property to infinitely many agents. Of course, every actual instance of the protocol will encompass only finitely many agents, but the proof provided here is uniform for any number of agents. The current work also serves as an example of reflecting meta-properties into specifications, allowing for the verification of additional properties.

Assuming a meta-level fairness property to prove liveness properties of the specification might be considered a circular argument, but without regress to meta-level fairness, no proof of the entrance property can be achieved. Keeping this in mind, our goal is to provide a reasonably simple and intuitive definition of fairness on the meta-level, that can be used for verification of the necessary properties, similar to any axiomatic approach where trust is based on simple axioms.

Future work we are foreseeing centers around the following points:

– Adaption of the methodology to other protocols: One probable candidate is the already mentioned Alternating Bit Protocol, where proofs for liveness properties are hard to obtain in other systems. We believe that a meta-specification similar to the one given here can be employed for this protocol, as well as others.
– Automatization of the method: Most of the steps done during lifting the specification to the meta-level are semi-automatic. It might be interesting to
provide built-in functionality to extend a given specification based on states with transition sequences.

– Characterization of strength of properties: We have seen that the mutual exclusion property can be proven by only looking at states, that the progress property only needs access to a state and its successor, but the entrance property needs access to all future states. We are working on a formal description of this concept called \( n \)-visibility.

We conclude with recapitulating the major contributions of this work: First of all, it provides an example for the inclusion of meta-concepts in a formal specification. Reflecting these meta-properties allows for the verification of additional properties, in particular liveness properties.

Furthermore, it is an example of a specification that spans all the corners of the CAFEOBJ cube, in particular mixing co-algebraic methods, infinite stream representation via hidden sorts, with transition sequence style modeling.

References

A Hybrid Method for the Verification and Synthesis of Parameterized Self-Stabilizing Protocols*

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Abstract. This paper presents a hybrid method for verification and synthesis of parameterized self-stabilizing protocols where algorithmic design and mechanical verification techniques/tools are used hand-in-hand. The core idea behind the proposed method includes the automated synthesis of self-stabilizing protocols in a limited scope (i.e., fixed number of processes) and the use of theorem proving methods for the generalization of the solutions produced by the synthesizer. Specifically, we use the Prototype Verification System (PVS) to mechanically verify an algorithm for the synthesis of weakly self-stabilizing protocols. Then, we reuse the proof of correctness of the synthesis algorithm to establish the correctness of the generalized versions of synthesized protocols for an arbitrary number of processes. We demonstrate the proposed approach in the context of an agreement and a coloring protocol on the ring topology.

Keywords: Mechanical Verification, Program Synthesis, Self-Stabilization, Parameterized Systems

1 Introduction

Self-stabilization is an important property of dependable distributed systems as it guarantees convergence in the presence of transient faults. That is, from any state/configuration, a Self-Stabilizing (SS) system recovers to a set of legitimate states (a.k.a. invariant) in a finite number of steps. Moreover, from its invariant, the executions of an SS system satisfy its specifications and remain in the invariant; i.e., closure. Nonetheless, design and verification of convergence are difficult tasks [10, 19] in part due to the requirements of (i) recovery from arbitrary states; (ii) recovery under distribution constraints, where processes can read/write only the state of their neighboring processes (a.k.a. their locality), and (iii) the non-interference of convergence with closure. Methods for algorithmic

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design of convergence [3, 4, 16, 13] can generate only the protocols that are correct up to a limited number of processes and small domains for variables. Thus, it is desirable to devise methods that enable automated design of parameterized SS systems, where a parameterized system includes several sets of symmetric processes that have a similar code up to variable re-naming.

Numerous approaches exist for mechanical verification of self-stabilizing systems most of which focus on synthesis and verification of specific protocols. For example, Qadeer and Shankar [32] present a mechanical proof of Dijkstra’s token ring protocol [10] in the Prototype Verification System (PVS) [34]. Kulkarni et al. [29] use PVS to mechanically prove the correctness of Dijkstra’s token ring protocol in a component-based fashion. Prasetya [31] mechanically proves the correctness of a self-stabilizing routing protocol in the HOL theorem prover [18]. Tsuchiya et al. [36] use symbolic model checking to verify several protocols such as mutual exclusion and leader election. Kulkarni et al. [28, 8] mechanically prove (in PVS) the correctness of algorithms for automated addition of fault tolerance; nonetheless, such algorithms are not tuned for the design of convergence. Most existing automated techniques [6, 27, 14, 16] for the design of fault tolerance enable the synthesis of non-parametric fault-tolerant systems. For example, Kulkarni and Arora [27] present a family of algorithms for automated design of fault tolerance in non-parametric systems, but they do not explicitly address self-stabilization. Abujarad and Kulkarni [4] present a method for algorithmic design of self-stabilization in locally-correctable protocols, where the local recovery of all processes ensures the global recovery of the entire distributed system. Farahat and Ebnenasir [13, 16] present algorithms for the design of self-stabilization in non-locally correctable systems. Jacobs and Bloem [22] show that, in general, synthesis of parameterized systems from temporal logic specifications is undecidable. They also present a semi-decision procedure for the synthesis of a specific class of parameterized systems in the absence of faults.

The contributions of this paper are two-fold: a hybrid method (Figure 1) for the synthesis of parameterized self-stabilizing systems and a reusable PVS theory for mechanical verification of self-stabilization. The proposed method includes a synthesis step and a theorem proving step. Our previous work [13, 16] enables the synthesis step where we take a non-stabilizing protocol and generate a self-stabilizing version thereof that is correct by construction up to a certain number of processes. This paper investigates the second step where we use the theorem prover PVS to prove (or disprove) the correctness of the synthesized protocol for an arbitrary number of processes; i.e., generalize the synthesized protocol. The synthesis algorithms in [13, 16] incorporate weak and strong convergence in existing network protocols; i.e., adding convergence. Weak (respectively, Strong) convergence requires that from every state there exists an execution that (respectively, every execution) reaches an invariant state in finite number of steps. To enable the second step, we first mechanically prove the correctness of the Add_Weak algorithm from [16] that adds weak convergence. As a result, any protocol generated by Add_Weak will be correct by construction. Moreover, the mechanical verification of Add_Weak provides a reusable theory in
PVS that enables us to verify the generalizability of small instances of different protocols generated by an implementation of Add Weak. If the mechanical verification succeeds, then it follows that the synthesized protocol is in fact correct for an arbitrary number of processes. Otherwise, we use the feedback of PVS to determine why the synthesized protocol cannot be generalized and re-generate a protocol that addresses the concerns reported by PVS. We continue this cycle of *synthesize and generalize* until we have a parameterized protocol. Notice that the theory developed in mechanical proof of Add Weak can also be reused for the mechanical verification of self-stabilizing protocols designed by means other than our synthesis algorithms. We demonstrate this reusability in the context of a coloring protocol (Section 6) and a binary agreement protocol (in [35]).

Fig. 1. A Hybrid method for the synthesis of parameterized self-stabilizing protocols.

**Organization.** Section 2 introduces basic concepts and presents their formal specifications in PVS. Then, Section 3 formally presents the problem of adding convergence to protocols. Sections 4 and 5 respectively present the specification and verification of Add Weak in PVS. Section 6 demonstrates the reusability and generalizability properties in the context of a graph coloring protocol. Section 7 discusses related work and Section 8 makes concluding remarks and presents future extensions of this work.

## 2 Formal Specifications of Basic Concepts

In this section, we define protocols, state predicates, computations and convergence, and present their formal specifications in PVS. The definitions of protocols and convergence are adapted respectively from [28] and [10].

### 2.1 Protocols

A protocol includes a set of processes, a set of variables and a set of transitions. Since we would like the specification of a protocol to be as general as possible, we impose little constraints on the notions of state, transitions, etc. Thus, the notations *state*, *variable*, and *domain* are all abstract and nonempty. Formally, we specify them by uninterpreted types state: Type+, Variable: Type+, Dom: Type+, where ‘+’ denotes the non-emptiness of the declared type. A *state predicate* is a set of states specified as StatePred: TYPE = set[state]. The concept of transition
is modeled as a tuple type of a pair of states Transition: Type = \([\text{state}, \text{state}]\) \[28\]. Likewise, an action is defined as a set of transitions, Action: Type = \(\text{set}[\text{Transition}]\). An action can be considered as an atomic guarded command “\(\text{grd} \rightarrow \text{stmt}\)”, where \(\text{grd}\) denotes a Boolean expression in terms of protocol variables and \(\text{stmt}\) is a set of statements that atomically update program variables when \(\text{grd}\) holds. We assume that Dom, Variable, and Action are finite types in our PVS specifications. A process is a tuple of a subset of variables that are readable by that process, a subset of variables that are writable by that process, and its subset of transitions. A protocol \(\text{prt}\) is a tuple of a finite set of processes, variables, and finite set of transitions. (In PVS, \(\text{"definedType}\colon \text{TYPE}^+ = [\]’’ declares the type \(\text{definedType}.\))

\[
\begin{align*}
\text{p_process} : \text{TYPE}^+ &= [\text{set[Variable]}, \text{set[Variable]}, \text{set[Transition]}] \\
\text{nd_Protocol} : \text{TYPE}^+ &= [\text{set[\text{p_process}]}, \text{set[Variable]}, \text{set[Transition]}]
\end{align*}
\]

The projection of a protocol \(\text{prt}\) on a state predicate \(I\), denoted \(\text{PrjOnS}(\text{prt}, I)\), includes the set of transitions of \(\text{prt}\) that start in \(I\) and end in \(I\). One can think of the projection of \(\text{prt}\) as a protocol that has the same set of processes and variables as those of \(\text{prt}\), but its transition set is a subset of \(\text{prt}\)’s transitions confined in \(I\). We model this concept by defining the following function, where \(Z\) is instantiated by transitions of \(\text{prt}\). (\(\text{proj}_k\) is a built-in function in PVS that returns the \(k\)-th element of a tuple.)

\[
\text{PrjOnS}(Z\colon \text{Action}, I\colon \text{StatePred})\colon \text{Action} = \\
\{t\colon \text{Transition} \mid t \in Z \land \text{proj}_1(t) \in I \land \text{proj}_2(t) \in I\}
\]

### 2.2 Distribution and Atomicity Models

We model the impact of distribution in a shared memory model by considering read and write restrictions for processes with respect to variables. Due to inability of a process \(P_j\) in reading some variables, each transition of \(P_j\) belongs to a group of transitions. For example, consider two processes \(P_0\) and \(P_1\) each having a Boolean variable that is not readable for the other process. That is, \(P_0\) (respectively, \(P_1\)) can read and write \(x_0\) (respectively, \(x_1\)), but cannot read \(x_1\) (respectively, \(x_0\)). Let \(\{x_0, x_1\}\) denote a state of this program. Now, if \(P_0\) writes \(x_0\) in a transition \((0, 0), (1, 0)\), then \(P_0\) has to consider the possibility of \(x_1\) being 1 when it updates \(x_0\) from 0 to 1. As such, executing an action in which the value of \(x_0\) is changed from 0 to 1 is captured by the fact that a group of two transitions \((0, 0), (1, 0)\) and \((0, 1), (1, 1)\) is included in \(P_0\).

In general, a transition is included in the set of transitions of a process \(j\) if and only if its associated group of transitions is included. Formally, any two transitions \((s_0, s_1)\) and \((s'_0, s'_1)\) in a group of transitions formed due to the read restrictions of a process \(P_j\) meet the following constraints, where \(r_j\) denotes the set of variables \(P_j\) can read: \(\forall v \in r_j : (v(s_0) = v(s'_0)) \land (v(s_1) = v(s'_1))\) and \(\forall v \notin r_j : (v(s_0) \neq v(s'_0)) \land (v(s_1) = v(s'_1))\), where \(v(s)\) denotes the value of a variable \(v\) in a state \(s\) that is represented by the \(\text{Val}(v, s)\) function in PVS. To enable the reusability of our PVS specifications, we specify our distribution model
as a set of axioms so one can mechanically prove convergence under different distribution and atomicity models.

In the following formal specifications, \( v \) is of type \( \text{Variable} \), \( p \) is of type \( \text{p_process} \), \( t \) and \( t' \) are of type \( \text{Transition} \), and \( \text{non_read} \) and \( \text{transition_group} \) are functions that respectively return the set of unreadable variables of the process \( p \) and the set of transitions that meet \( \forall v : v \in r_j : (v(s_0) = v(s_0')) \land (v(s_1) = v(s_1')) \) for a transition \( t = (s_0, s_1) \) and its groupmate \( t' = (s_0', s_1') \).

\[
\text{AXIOM subset?(proj_2(p),proj_1(p))} \quad /\!/ \text{Writable variables are a subset of readable variables.}
\]

\[
\text{AXIOM member(t',transition_group(p,t,prt)) AND member(v,Non_read(p,prt)) IMPLIES Val(v,proj_1(t)) = Val(v,proj_2(t)) AND Val(v,proj_1(t')) = Val(v,proj_2(t'))}
\]

\( \text{member}(x,X) \) and \( \text{subset?}(X,Y) \) respectively represent the membership and subset predicates in a set-theoretic sense.

### 2.3 Computation

A \textit{computation} of a protocol \( \text{prt} \) is a sequence \( A \) of states, where \( (A(i), A(i+1)) \) represents a transition of \( \text{prt} \) executed by some action of \( \text{prt} \). In a more general term, a computation of any set of transitions \( Z \), is a sequence of states in which every state can be reached from its predecessor by a transition in \( Z \). Thus, we define the following function to return the set of computations generated by the set of transitions \( Z \).

\[
\text{COMPUTATION}(Z: \text{Action}) : \text{set}[\text{sequence}[\text{state}]] = \{A: \text{sequence}[\text{state}] | \forall (n : \text{nat}) : ((A(n), A(n + 1)) \in Z)\}
\]

A \textit{computation prefix} of a protocol \( \text{prt} \) is a \textit{finite} sequence of states where each state is reached from its predecessor by a transition of \( \text{prt} \). Kulkarni et al. [28] specify a prefix as an infinite sequence in which only a finite number of states are used. By contrast, we specify a computation prefix as a finite sequence type. We believe that it is more natural and more accurate to model the concept of prefix by finite sequences. Our experience also shows that modeling computation prefixes as finite sequences simplifies formal specification and verification of reachability and convergence while saving us several definitions that were required in [28] to capture the length of the prefix. We use the predicate \( \text{Condi_prefix?}(A,Z) \) that holds when all transitions \( (A(i), A(i + 1)) \) of a sequence \( A \) belong to a set of transitions \( Z \). The notation \( A'\text{length} \) denotes the length of the sequence \( A \), \( A'\text{seq}(i) \) returns the \( i \)th element of sequence \( A \) and \( \text{below}[k] \) represents natural values less than \( k \). The function \( \text{PREFIX} \) returns the set of computation prefixes generated by transitions of \( Z \).

\[
\text{Pos_F_S: TYPE = \{c:finite_sequence[\text{state}] | c'\text{length} > 0\}}
\]

\[
\text{Condi_Prefix?(A:Pos_F_S,Z:Action):bool= FORALL(i: below[A'length-1] ) : member((A'\text{seq}(i), A'\text{seq}(i+1)), Z)}
\]

\[
\text{PREFIX(Z: Action) : set[Pos_F_S]= \{A:Pos_F_S | Condi_Prefix?(A,Z) \}}
\]
2.4 Closure and Convergence

A state predicate $I$ is closed in a protocol $prt$ iff every transition of $prt$ that starts in $I$ also terminates in $I$ \cite{19,5}. The closed predicate checks whether a set of transitions $Z$ is actually closed in a state predicate $I$.

$$\text{closed?}(I: \text{StatePred}, Z: \text{Action}) = \text{FORALL}(t: \text{Transition} \mid (\text{member}(t,Z)) \text{ AND } \text{member}((\text{proj}_1(t), I)) : \text{member}((\text{proj}_2(t), I))$$

A protocol $prt$ weakly converges to a non-empty state predicate $I$ iff from every state $s$, there exists at least one computation prefix that reaches some state in $I$ \cite{19,5}. A strongly converging protocol guarantees that every computation from $s$ will reach some state in $I$. Notice that any strongly converging protocol is also weakly converging, but the reverse is not necessarily true. A protocol $prt$ is weakly (respectively, strongly) self-stabilizing to a state predicate $I$ iff (1) $I$ is closed in $prt$, and (2) $prt$ weakly (respectively, strongly) converges to $I$.

3 Problem Statement

The problem of adding convergence (from \cite{16}) is a transformation problem that takes as its input a protocol $prt$ and a state predicate $I$ that is closed in $prt$. The output of Problem 1 is a revised version of $prt$, denoted $prt_{ss}$, that converges to $I$ from any state. Starting from a state in $I$, $prt_{ss}$ generates the same computations as those of $prt$; i.e., $prt_{ss}$ behaves similar to $prt$ in $I$.

**Problem 1. Add Convergence**

- **Input:** (1) A protocol $prt$; (2) A state predicate $I$ such that $I$ is closed in $prt$; and (3) A property of $L_s$ converging, where $L_s \in \{\text{weakly, strongly}\}$.
- **Output:** A protocol $prt_{ss}$ such that: (1) $I$ is unchanged; (2) the projection of $prt_{ss}$ on $I$ is equal to the projection of $prt$ on $I$, and (3) $prt_{ss}$ is $L_s$ converging to $I$. Since $I$ is closed in $prt_{ss}$, it follows that $prt_{ss}$ is $L_s$ self-stabilizing to $I$.

Previous work \cite{19,16} shows that weak convergence can be added in polynomial time (in the size of the state space), whereas adding strong convergence is known to be an NP-complete problem \cite{26}. Farahat and Ebnenasir \cite{16,13} present a sound and complete algorithm for the addition of weak convergence and a set of heuristics for efficient addition of strong convergence. While one of our objectives is to develop a reusable proof library (in PVS) for mechanical verification of both weak and strong convergence, the focus of this paper is mainly on enabling the mechanical verification of weak convergence for parameterized systems. Algorithm 1 provides an informal and self-explanatory representation of the $\text{Add\_Weak}$ algorithm presented in \cite{16}.

Mechanical verification of the soundness of $\text{Add\_Weak}$ ensures that any protocol synthesized by $\text{Add\_Weak}$ is correct by construction. Moreover, the lemmas and theorems developed in mechanical verification of $\text{Add\_Weak}$ provide a reusable framework for mechanical verification of different protocols that we generate using our synthesis tools \cite{16,25}. The verification of synthesized protocols increases our confidence in the correctness of the implementation of $\text{Add\_Weak}$ and helps us to generalize small instances of weakly converging protocols to their parameterized versions.
Algorithm 1: Add_Weak

Input: $prt$ : nd_protocol, $I$ : statePred
Output: set[Transition]: // Set of transitions of a weakly self-stabilizing version of $prt$

1: Let $Δ_{\text{prt}}$ be the set of transition groups of $prt$.
2: Let $Δ_{\text{converge}}$ be the set of transition groups that adhere to read/write restrictions of processes of $prt$, but exclude any transition starting in $I$;
3: $Δ_{\text{ws}} = Δ_{\text{prt}} \cup Δ_{\text{converge}}$;
4: $\text{noPrefix} := \{ s : \text{state} \mid (s \notin I) \land \text{(there is no computation prefix using transitions of } Δ_{\text{ws}} \text{ that can reach a state in } I)\}$
5: If ($\text{noPrefix} \neq \emptyset$) then weak convergence cannot be added to $prt$; return;
6: return $Δ_{\text{ws}}$;

4 Specification of Add_Weak

This section presents the highlights of the formal specification of Add_Weak in PVS. (The complete PVS specifications are available at http://asd.cs.mtu.edu/projects/mechVerif/ss.html.) We start by specifying the basic components used in the Add_Weak algorithm, namely the transition predicates $Δ_{\text{prt}}, Δ_{\text{converge}}$ and $Δ_{\text{ws}}$, and the state predicate $\text{noPrefix}$.

Notation. In the subsequent formal specifications, we use the identifiers $\text{Delta}_\text{prt}, \text{Delta}_\text{Converge}$ and $\text{Delta}_\text{ws}$ corresponding to the variables $Δ_{\text{prt}}, Δ_{\text{converge}}$ and $Δ_{\text{ws}}$ in Add_Weak. The function $\text{transition groups proc}(p,prt)$ returns the set of transition groups of a process $p$ of a protocol $prt$.

\[
\text{Delta}_\text{Converge}(prt: \text{nd_protocol}, I: \text{StatePred}) : \text{set[set[Transition]]} = \{ gg : \text{set[Transition]} \mid \exists (p, p' : \text{process}) : \text{member}(p, \text{proj}_1(prt)) \land \text{member}(gg, \text{transition groups proc}(p, prt)) \land \forall (t : \text{Transition}) : \text{member}(t, gg) \land \text{member}((\text{proj}_1(t), I), gg) \}
\]

We find it useful to define a dependent type of all prefixes $A$ of a set of transitions $Z$ and we call it $\text{PREFIX}_T(Z)$. Furthermore, we formally specify the concept of reachability as follows:

\[
\text{Reach from?}(Z : \text{Action}, A : \text{PREFIX}_T(Z), s_0 : \text{state}, I : \text{StatePred}) : \text{bool} = \exists (j : \text{below}[A'\text{length}]) : A'\text{seq}(0) = s_0 \land \text{member}(A'\text{seq}(j), I)
\]

The predicate $\text{Reach from}$ returns true iff a state predicate $I$ is reachable from a state $s_0$ using computation prefixes of $Z$. To specify the set of states $\text{noPrefix}$, we first specify a predicate $\text{condi_noPrefix}$ that determines if for a protocol $prt$, a state predicate $I$ and a state $s_0$, no state in $I$ can be reached from $s_0$ by computation prefixes of $prt$.

\[
\text{condi_noPrefix?}(prt: \text{nd_protocol}, I: \text{StatePred}, s_0 : \text{state}) : \text{bool} = \forall (g : \text{Action}, A : \text{PREFIX}_T(g) : \text{member}(g, \text{Delta}_\text{ws}(prt, I)) \land \text{member}(A, \text{PREFIX}(g)) \land A'\text{seq}(0) = s_0) : \text{NOT } (\text{Reach from?}(g, A, s_0, I))
\]

We then specify the state predicate $\text{noPrefix}$ in Add_Weak as follows:
no_Prefix(prt:nd.Protocol,I:StatePred):set[state] = {s0:state | NOT(member(s0,I)) AND cond_no_Prefix?(prt,I,s0)}

We also specify Add_Weak as a function that returns a set of transitions.

Add_weak(prt:nd.Protocol, I:{II:StatePred | closed?(II,proj 3(prt))}):set[Transition] = COND
empty?(no_Prefix(prt,I)) – > Delta_ws(prt,I), ELSE – > proj 3(prt) ENDCOND

5 Verification of Add_Weak

In order to prove the soundness of Add_Weak, we check if (1) I is unchanged; (2) the projection of ∆ws on I is equal to the projection of ∆prt on I, and (3) ∆ws is weakly converging to I. The first constraint holds trivially since no step of Add_Weak adds/removes a state to/from I. Next, we present a set of lemmas and theorems that prove the other two constraints of Problem 1.

5.1 Verifying the Equality of Projections on Invariant

In this section, we prove that Constraint 2 of Problem 1 holds for the output of Add_Weak, denoted by a protocol whose set of transitions is ∆ws. Our proof obligation is to show that the projection of ∆ws on I is equal to the projection of ∆prt on I. We decompose this into two set inclusion obligations of PrjOnS(∆prt,I) ⊆ PrjOnS(∆ws,I) and PrjOnS(∆ws,I) ⊆ PrjOnS(∆prt,I). Notice that, by assumption, closed?(I,∆prt) is true.

Lemma 1. PrjOnS(∆prt,I) is a subset of PrjOnS(∆ws,I).

Proof. The proof is straightforward since by construction we have ∆ws = ∆prt ∪ ∆converge.

Lemma 2. PrjOnS(∆ws,I) is a subset of PrjOnS(∆prt,I).

Proof. If a transition t = (s0,s1) is in PrjOnS(∆ws,I) then s0 ∈ I. Since ∆ws = ∆prt ∪ ∆converge, either t ∈ ∆prt or t ∈ ∆converge. By construction, ∆converge excludes any transition starting in I including t. Thus, t must be in ∆prt. Since s0 ∈ I, it follows that t ∈ PrjOnS(∆prt,I).


5.2 Verifying Weak Convergence

In this section, we prove the weak convergence property (i.e., Constraint 3 of Problem 1) of the output of Add_Weak. Specifically, we show that from any state s0 ∈ ¬I, there is a prefix A in PREFIX(∆ws) such that A reaches some state in I. Again, we observe that, an underlying assumption in this section is that closed?(I,∆prt) holds. For a protocol prt and a predicate I that is closed in prt and a state s /∈ I, we have:
Lemma 3. If empty?(no Prefix(prt,I)) holds then condi_no_Prefix?(prt,I,s) returns false.

Lemma 4. If condi_no_Prefix?(prt,I,s) returns false for some s /∈ I, then there exists a sequence of states A and a set of transitions Z such that Z ∈ Δ_ws, A ∈ PREFIX(Z), A(0)=s holds, and Reach_from?(Z,A,s,I) returns true.

Lemma 4 implies that when Add Weak returns, the revised version of prt guarantees that there exists a computation prefix to I from any state outside I; hence weak convergence. This is due to the fact that A is a prefix of Δ_ws.

Theorem 2. If empty?(no Prefix(prt,I)) holds and s /∈ I then there exists a sequence of states A that starts from s and A ∈ PREFIX(Δ_ws) and Reach_from?(Δ_ws(prt,I),A,s,I) returns true. (PVS specifications and proofs are available at http://asd.cs.mtu.edu/projects/mechVerif/ss.html.)

6 Reusability and Generalizability

In this section, we demonstrate how the lemmas and theorems proved for the soundness of Add Weak can be reused in proving the correctness of a graph coloring protocol and in generalizing it. Due to space constraints, we omit the proof of correctness of a binary agreement protocol (which is available in [35]). Reusability enables us to instantiate the abstract concepts/types (e.g., state predicate I and actions of a protocol) for a concrete protocol and reuse the mechanical proof of Add Weak to prove the weak convergence of that protocol. Generalizability determines whether a small instance of a protocol synthesized by our implementation of Add Weak [16] can be proven to be correct for an arbitrary number of processes.

TR(m,n): Coloring on a ring of n processes with m > 2 colors. We have used the Stabilization Synthesizer (STSyn) [16] tool to automatically generate the 3-coloring protocol for rings of up to 40 processes (i.e., n < 41). Nonetheless, due to scalability issues, STSyn cannot synthesize a self-stabilizing 3-coloring protocol for n > 40. In this section, we apply the proposed approach of synthesize in small scale and generalize to prove (or disprove) that the synthesized 3-coloring protocol is correct for rings of size greater than 40 and with more than 2 colors (i.e., m > 2).

The coloring protocol, denoted TR(m,n), includes n > 3 processes located along a bidirectional ring. Each process Pj has a local variable cj with a domain of m > 2 values representing m colors. Thus, the set of variables of TR(m,n) is V_TR(m,n) = {c0, c1, ..., cn-1}. Each process Pj can read {cj⊖1, cj, cj⊕1}, and is allowed to write only cj, where ⊖ and ⊕ denote addition and subtraction modulo n respectively. The set of legitimate states of TR(m,n) includes the states where no two neighboring processes have the same color. Formally, I_coloring = ∀j : 0 ≤ j < n : cj ≠ cj⊖1. The coloring protocol has applications in several domains such as scheduling, bandwidth allocation, register allocation, etc. It is known that if m > d, where d is the max degree in the topology graph of the system, then
the coloring problem is solvable. For this reason, we have \( m > 2 \) for the ring.

Using STSyn [16], we have automatically generated the following action for each process \( P_j \) (\( 0 \leq j < 41 \)):

\[
A_j : (c_j = c_{j\oplus 1}) \lor (c_j = c_{j\ominus 1}) \rightarrow c_j := \text{other}(c_{j\oplus 1}, c_{j\ominus 1})
\]  

(1)

If \( P_j \) has the same color as that of one of its neighbors, then \( P_j \) uses the function \( \text{other}(c_{j\oplus 1}, c_{j\ominus 1}) \) to non-deterministically set \( c_j \) to a color different from \( c_{j\ominus 1} \) and \( c_{j\oplus 1} \). While \( TR(3,n) \) is correct by construction for \( n \leq 40 \), we would like to investigate whether \( TR(3,n) \) is weakly stabilizing for \( n > 40 \). The reuse of mechanical proof of Add\_Weak greatly simplifies the proof of generalization of the synthesized protocol.

6.1 PVS Specification of Coloring

This section presents the PVS specification of \( TR(m,n) \). First, we instantiate the basic types in the PVS specification of Add\_Weak for the coloring protocol. Then, we present the specifications of some functions that we use to simplify the verification tasks. Finally, we specify the processes and the protocol itself.

**Basic types.** We first define the parameterized type \( \text{COLORS}: \text{below}[m] \) to capture the colors and the size of variable domains. Then, we define a state as a finite sequence of colors of length \( n \); i.e., \( \text{STC: NONEMPTY_TYPE \{ s: \text{finseq } | \text{s.length=n} \} } \). Since each variable \( c_j \) holds two pieces of information namely the process position in the ring and the color, we model their type by the tuple type \( \text{ndx_varb:TYPE=}=[\text{COLORS, below}[n]] \). The predicate \( \text{is nbr?}(K:\text{ndx_varb}, L:\text{ndx_varb}) \) returns true iff \( K \) and \( L \) are two neighboring processes; i.e., \( \text{mod(abs(K'2-L'2),n)} \leq 1 \), where \( K'2 \) denotes the second element of the pair \( K \) (which is the position of \( K \) in the ring). Likewise, we define the predicate \( \text{is bad
br?}(K:\text{ndx_varb}, L:\text{ndx_varb}) \) that holds iff \( \text{is nbr?}(K,L) \) holds and \( K'1 = L'1 \).

To capture the *locality* of process \( j \), we define the non-empty dependent type \( \text{nbr\_v}(K:\text{ndx_varb}): \text{TYPE=} = \{L: \text{ndx_varb } | \text{is_nbr?}(K,L) \} \). Likewise, we define the type \( \text{bad
br\_v}(K:\text{ndx_varb}): \text{TYPE=} = \{L: \text{ndx_varb } | \text{is bad
br?}(K,L) \} \) to capture the set of neighbors of a process that have the same color as that process. The function \( \text{nbr\_colors}(K:\text{ndx_varb}): \text{set}[\text{COLORS}] \) returns the set of colors of the immediate neighbors of a process.

**Functions.** In order to simplify the verification of convergence (in Section 6.2), we associate the subsequent functions with a global state. For example, we define a function \( \text{ValPos}(s:\text{STC},j: \text{below}[n]): \text{ndx_varb} = (s'\text{seq}(j), j) \) that returns the value and the position of process \( j \) in a global state \( s \) as a tuple of type \( \text{ndx_varb} \). An example use of this function is \( \text{Val}(s:\text{STC}, \text{L:ndx_varb}) = \text{ValPos}(s, L'2) \). Moreover, the predicate \( \text{nbr\_is\_bad?}(s:\text{STC},j: \text{below}[n]): \text{bool} = \text{nonempty?}(\text{bad
br\_v}(\text{ValPos}(s,j))) \) returns true iff for an arbitrary state \( s \) and a process \( j \) the set of bad neighbors of the variable \( \text{ValPos}(s,j) \) is nonempty; we refer to such a case by saying \( s \) is *corrupted at process\( j \). Notice that an illegitimate state can be corrupted at more than one position. We also define the predicate \( \text{nbr\_is\_good?}(s:\text{STC},j: \text{below}[n]): \text{bool} \) as the negation of the predicate \( \text{nbr\_is\_bad?}(s:\text{STC},j: \text{below}[n]): \text{bool} \). We also define
the predicate is_LEGT?(s:STC):bool= Forall (j:below[n]): nbr_is_good?(s,j) that returns true if s is a legitimate state. Thus, the set of illegitimate states is specified as $S_{ill}: TYPE= \{ s:STC \mid \text{not is_LEGT?(s)} \}$. The association of a global state to functions and types enables us to import the PVS theory of Add_Weak with the following types $[STC, below[m], ndx\_varb, [ndx\_varb, STC \rightarrow below[m]]]$ to reuse its already defined types and functions in specifying $TR(m, n)$ as follows.

**Specification of a process of TR(m,n).** For an arbitrary global state $s$ and a process $j$, we define the function $READ_p$ which returns all readable variables of process $j$.

$$READ_p(s:STC, j:below[n]): set[ndx\_varb] = \{ L:nbr = ValPos(s,j) \}.$$  

Similarly, we define the function $WRITE_p$ which returns the variables that process $j$ can write.

$$WRITE_p(s:STC, j:below[n]): set[ndx\_varb] = \{ L:ndx\_varb = ValPos(s,j) \}$$

We now define the function $DELTA_p$ that returns the set of transitions belonging to process $j$ if process $j$ is corrupted in the global state $s$; i.e., $j$ has a bad neighbor.

$$DELTA_p(s:STC, j:below[n]): set[Transition] = \{ tr: Transition = (s,action(s,j,ValPos(s,j),c)) \}$$

The function $action(s,j,ValPos(s,j),c)$ returns the state reached when process $j$ acts to correct its corrupted state. Formally, we define $action(s,j,ValPos(s,j),c)$ as follows: (The LAMBDA abstractions in PVS enable us to specify binding expressions similar to quantified statements in predicate logic.)

$$action(s:STC,j:below[n],K:ndx\_varb, C:bad\_nbr\_v(ValPos(s,j))): STC = (# length := n, seq := (LAMBDA (i:below[n]):IF \ i = j \ THEN o\_ther(ValPos(s,j)) \ ELSE s(i) ENDIF) #)$$

We specify the function $o\_ther$ to randomly choose a new color other than the corrupted one. To this end, we use the $epsilon$ function over the full set of colors minus the set of colors of the neighbors of the corrupted process. Formally, we have $other(K,ndx\_varb):COLORS = epsilon(difference(fullset\_colors,nbr\_colors(K)))$, where $fullset\_colors:set[COLORS] = \{ cl:COLORS \mid \text{TRUE} \}$. Thus, the specification of a process of the protocol $TR(m, n)$ is as follows:

$$PRS_p(s:STC, j:below[n]): p\_process = (READ_p(s,j), WRITE_p(s,j), DELTA_p(s,j))$$

The parameterized specification of the $TR(m,n)$ protocol. We define the $TR(m,n)$ protocol as the type $TR_m_n(s:STC): nd\_Protocol = (PROC\_prt(s),VARB\_prt(s),DELTA\_prt(s))$, where the parameters are defined as follows:

$$PROC\_prt(s:STC): set[p\_process]=\{ p:p\_process \mid \text{Exists (j:below[n]): p = PRS_p(s,j)} \}$$

$$VARB\_prt(s:STC): set[ndx\_varb]=\{ v:ndx\_varb \mid \text{Exists (j:below[n]): member(v, WRITE_p(s,j))} \}$$

$$Delta\_prt(s:STC): set[Transition]=\{ tr:Transition \mid \text{Exists (j:below[n]): member(tr,DELTA_p(s,j))} \}$$

**6.2 Mechanical Verification of Parameterized Coloring**

We now prove the weak convergence of $TR(m, n)$ for $m > 2$ and $n > 40$. To this end, we show that the set no_Prefix of $TR(m, n)$ is actually empty. The
proof of emptiness of no_Prefix is based on a prefix constructor function that demonstrates the existence of a computation prefix $\sigma$ of $TR(m, n)$ from any arbitrary illegitimate state $s$ such that $\sigma$ includes a state in $I$. Subsequently, we instantiate Theorem 2 for $TR(m, n)$. Due to space constraints, we omit the details of the mechanical proof (which are available in [35]).

Theorem 3. Let $prt$ be $TR(m, n)$. If $\text{closed?}(I, \text{no_Prefix}(prt, I))$ and $\text{empty?}(\text{no_Prefix}(prt, I))$ hold and $s \notin I$ then there exists a sequence of states $A$ that starts from $s$ and $A \in \text{PREFIX}(\text{Delta_wis}(prt, I))$ and $\text{Reach_from?}(A, s, I)$ returns true.

7 Discussion and Related Work

This section discusses the impact of the proposed approach and the related work. Self-stabilization is an important property for networked systems, be it a network-on-chip system or the Internet. There are both hardware [11] and software systems [12] that benefit from the resilience provided by self-stabilization. Thus, it is important to have an abstract specification of self-stabilization that is independent from hardware or software. While several researchers [32, 29] have utilized theorem proving to formally specify and verify the self-stabilization of specific protocols, this paper presents a problem-independent specification of weak convergence that enables potential reuse of efforts in the verification of convergence of different protocols.

One of the fundamental impediments before automated synthesis of self-stabilizing protocols from their non-stabilizing versions is the scalability problem. While there are techniques for parameterized synthesis [22, 24] of concurrent systems, such methods are not directly useful for the synthesis of self-stabilization due to several factors. First, such methods are mostly geared towards synthesizing concurrent systems from formal specifications in some variant of temporal logic. Second, in the existing parameterized synthesis methods the formal specifications are often parameterized in terms of local liveness properties of individual components (e.g., progress for each process), whereas convergence is a global liveness property. Third, existing methods often consider the synthesis from a set of initial states that is a proper subset of the state space rather than the entire state space itself (which is the case for self-stabilization). With this motivation, our contributions in this paper enable a hybrid method based on synthesis and theorem proving that enables the generalization of small instances of self-stabilizing protocols generated by our tools [25].

Related work. Kulkarni and Bonakdarpour’s work [28, 8] is the closest to the proposed approach in this paper. As such, we would like to highlight some differences between their contributions and ours. First, in [28], the authors focus on mechanical verification of algorithms for the addition of fault tolerance to concurrent systems in a high atomicity model where each process can read and write all system variables in one atomic step. One of the fault tolerance requirements they consider is nonmasking fault-tolerance, where a nonmasking system guarantees recovery to a set of legitimate states from states reachable by faults and not necessarily from the entire state space. Moreover, in [8], Kulkarni and
Bonakdarpour investigate the mechanical verification of algorithms for the addition of multiple levels of fault tolerance in the high atomicity model. In this paper, our focus is on self-stabilization in distributed systems where recovery should be provided from any state and high atomicity actions are not feasible.

Methods for the verification of parameterized systems can be classified into four major approaches, which do not directly address SS systems. Abstraction techniques [21, 30, 15] generate a finite-state model of a parameterized system and then reduce the verification of the parameterized system to the verification of its finite model. Network invariant approaches [37, 23, 20] find a process that satisfies the property of interest and is invariant to parallel composition. Logic program transformations and inductive verification methods [33, 17] encode the verification of a parameterized system as a constraint logic program and reduce the verification of the parameterized system to the equivalence of goals in the logic program. In regular model checking [9, 2], system states are represented by grammars over strings of arbitrary length, and a protocol is represented by a transducer. Abdulla et al. [1] also investigate reachability of unsafe states in symmetric timed networks and prove that it is undecidable to detect livelocks in such networks. Bertrand and Fournier [7] also focus on the verification of safety properties for parameterized systems with probabilistic timed processes.

8 Conclusion and Future Work

This paper focuses on exploiting theorem proving for the generalization of synthesized self-stabilizing protocols that are correct in a finite scope (i.e., up to a small number of processes). We are particularly interested in weak stabilization where reachability to legitimate states is guaranteed from any state. The contributions of this paper comprise a component of a hybrid method for verification and synthesis of parameterized self-stabilizing network protocols (see Figure 1). This paper specifically presents a mechanical proof for the correctness of the Add_Weak algorithm from [16] that synthesizes weak convergence. This mechanical proof provides a reusable theory in PVS for the proof of weakly stabilizing systems in general (irrespective of how they have been designed). The success of mechanical proof for a small synthesized protocol shows the generality of the synthesized solution for arbitrary number of processes. We have demonstrated the proposed approach in the context of a binary agreement protocol (in [35]) and a graph coloring protocol (Section 6).

We will extend this work by reusing the existing PVS theory for mechanical proof of algorithms (in [16]) that design strong convergence. Moreover, we are currently investigating the generalization of more complicated protocols (e.g., leader election, maximal matching, consensus) using the proposed approach.

References


Drill & Join
A method for exact inductive program synthesis

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Abstract. In this paper we propose a novel semi-supervised active machine-learning method, based on two recursive higher-order functions that can inductively synthesize a functional computer program. Based on properties formulated using abstract algebra terms, the method uses two combined strategies: to reduce the dimensionality of the Boolean algebra where a target function lies and to combine known operations belonging to the algebra, using them as a basis to build a program that emulates the target function. The method queries for data on specific points of the problem input space and build a program that exactly fits the data. Applications of this method include all sorts of systems based on bitwise operations. Any functional computer program can be emulated using this approach. Combinatorial circuit design, model acquisition from sensor data, reverse engineering of existing computer programs are all fields where the proposed method can be useful.

1 Introduction

Induction means reasoning from specific to general. In the case of inductive learning from examples, the general rules are derived from input/output (I/O) examples or answers from questions. Inductive machine learning has been successfully applied to a variety of classification and prediction problems [1] [3].

Inductive program synthesis (IPS) builds from examples the computation required to solve a problem. The problem must be formulated as a task of learning a concept from examples, referred to as inductive concept learning [14]. A computer program is automatically created from an incomplete specification of the concept to be implemented, also referred as the target function [3].

Research on inductive program synthesis started in the seventies. Since then it has been studied in several different research fields and communities such as artificial intelligence (AI), machine learning, inductive logic programming (ILP), genetic programming, and functional programming [3] [8].

One basic approach to IPS is to simply enumerate programs of a defined set until one is found which is consistent with the examples. Due to combinatorial explosion, this general enumerative approach is too expensive for practical use. Summers [9] proposed an analytical approach to induce functional Lisp programs without search. However, due to strong constraints imposed on the forms of I/O-examples and inducible programs in order to avoid search, only relatively simple functions can be induced. Several variants and extensions of Summers’ method
have been proposed, like in [7]. An overview is given in [10]. Kitzelmann [6] proposed a combined analytical and search-based approach. Albarghouthi [5] proposed ESCHER, a generic algorithm that interacts with the user via I/O examples, and synthesizes recursive programs implementing intended behavior.

Hybrid methods propose the integration of inductive inference and deductive reasoning. Deductive reasoning usually requires high-level specifications as a formula in a suitable logic, a background theory for semantic correctness specification, constraints or a set of existing components as candidate implementations. Some examples of hybrid methods are the syntax-guided synthesis [13], the sciduction methodology [19] and the oracle-guided component-based program synthesis approach [18].

IPS is usually associated to functional programming [8]. In functional code the output value of a function depends only on the arguments that are input to the function. It is a declarative programming paradigm. Side effects that cause change in state that do not depend on the function inputs are not allowed. Programming in a functional style can usually be accomplished in languages that aren’t specifically designed for functional programming. In early debates around programming paradigms, conventional imperative programming and functional programming were compared and discussed. John Backus [16], in his work on programs as mathematical objects, supported the functional style of programming as an alternative to the “ever growing, fat and weak” conventional programming languages. Backus identified as inherent defects of imperative programming languages their inability to effectively use powerful combining forms for building new programs from existing ones, and their lack of useful mathematical properties for reasoning about programs. Functional programming, and more particularly function-level programming, is founded on the use of combining forms for creating programs that allow an algebra of programs.

Our work is distinguishable from previous works on IPS in a number of ways:

- Our method is based on function-level programming. Programs are built directly from programs that are given at the outset, by combining them with program-forming operations.
- Our method is based on active learning. Active learning is a special case of machine learning in which the learning algorithm can control the selection of examples that it generalizes from and can query one or more oracles to obtain examples. The oracles could be implemented by evaluation/execution of a model on a concrete input or they could be human users. Most existing IPS methods supply the set of examples at the beginning of the learning process, without reference to the learning algorithm. Some hybrid synthesis methods use active learning, as in [13] [19] and [18], to generate examples to a deductive procedure. Our method defines a learning protocol that queries the oracle to obtain the desired outputs at new data points.
- Our method generates programs on a very low-level declarative language, compatible with most high-level programming languages. Most existing methods are conceived considering and restricted to specific high-level source languages. Our whole method is based on Boolean algebra. Inputs and out-
puts are bit vectors and the generated programs are Boolean expressions. The synthesized program can also be used for a combinatorial circuit design describing the sequences of gates required to emulate the target function. Research on reconfigurable supercomputing is very interested in providing compilers that translate algorithms directly into circuit design expressed in an hardware description language. They want to avoid the high cost of having to hand-code custom circuit designs [15].

- The use of Boolean algebra and abstract algebra concepts at the basis of the method defines a rich formalism. The space where a program is to be searched, or synthesized, corresponds to a well defined family of operations that can be reused, combined and ordered.

- Our method can be applied to general purpose computing. A program generated using our method, computing an output bit vector from an input bit vector, is equivalent to a system of Boolean equations or a set of truth tables. However, bit vectors can also be used to represent any kind of complex data types, like floating point numbers and text strings. A functionally complete set of operations performed on arbitrary bits is enough to compute any computable value. In principle, any Boolean function can be built-up from a functionally complete set of logic operators. In logic, a functionally complete set of logical connectives or Boolean operators is one which can be used to express all possible truth tables by combining members of the set into a Boolean expression. Our method synthesizes Boolean expressions based on the logic operators set \{\text{XOR}, \text{AND}\} which is functionally complete.

- If the problem has a total functional behavior and enough data is supplied during the learning process our method can synthesize the exact solution.

This paper is organized as follows. Sections 2 and 3 review some relevant mathematical concepts. Sections 4 and 5 describe the mathematics of the method. Section 6 presents the method itself and how the programs are synthesized. Section 7 presents the main algorithms. Section 8 shows a Common Lisp implementation of the simplest version of the method. Sections 9, 10 and 11 close the document discussing the method, possible applications and future work.

2 Boolean ring, F2 field, Boolean polynomials and Boolean functions

A Boolean ring is essentially equivalent to a Boolean algebra, with ring multiplication corresponding to conjunction (\&\&) and ring addition to exclusive disjunction or symmetric difference (\oplus or XOR). In logic, the combination of operators \oplus (XOR or exclusive OR) and \&\& (AND) over elements true, false produce the Galois field F2 which is extensively used in digital logic and circuitry [12]. This field is functionally complete and can represent any logic obtainable with the system (\&\&\|) and can also be used as a standard algebra over the set of the in-
tegers modulo 2 (binary numbers 0 and 1). Addition has an identity element (false) and an inverse for every element. Multiplication has an identity element (true) and an inverse for every element but false.

Let $B$ be a Boolean algebra and consider the associated Boolean ring. A Boolean polynomial in $B$ is a string that results from a finite number of Boolean operations on a finite number of elements in $B$. A multivariate polynomial over a ring has a unique representation as a xor-sum of monomials. This gives a normal form for Boolean polynomials:

$$
\bigoplus_{J \subseteq \{1,2,\ldots,n\}} a_J \prod_{j \in J} x_j
$$

where $a_J \in B$ are uniquely determined. This representation is called the algebraic normal form. A Boolean function of $n$ variables $f : \mathbb{Z}_2^n \rightarrow \mathbb{Z}_2$ can be associated with a Boolean polynomial by deriving an algebraic normal form.

3 Abstract algebra and higher order functions

Let us consider a generic functional setting having as domain the set of bit strings of a finite, defined, length and as range the set $\{true, false\}$ or the binary numbers 0 and 1 represented by one bit. This setting can represent the inputs and output of a logic proposition, a Boolean function, a truth table or a fraction of a functional program corresponding to one of its output bits.

In abstract algebra terms, this setting will define a finitary Boolean algebra consisting of a finite family of operations on $\{0,1\}$ having the input bit string as their arguments. The length of the bit string will be the arity of the operations in the family. An $n$-ary operation can be applied to any of $2^n$ possible values of its $n$ arguments. For each choice of arguments an operation may return 0 or 1, whence there are $2^{2^n}$ different operations in the family. In this functional setting, IPS could be seen as the synthesis of an operation that fits a set of I/O examples inside its family. The Boolean algebra defines our program space.

Let $V^n$ be the set of all binary words of length $n$, $|V^n| = 2^n$. The Boolean algebra $B$ on $V^n$ is a vector space over $\mathbb{Z}_2$. Because it has $2^{2^n}$ elements, it is of dimension $2^n$ over $\mathbb{Z}_2$. This correspondence between an algebra and our program space defines some useful properties. The operations in a family need not be all explicitly stated. A basis is any set of operators from which the remaining operations can be obtained by composition. A Boolean algebra may be defined from any of several different bases. To be a basis is to yield all other operations by composition, whence any two bases must be intertranslatable.

A basis is a linearly independent spanning set. Let $v_1, \ldots, v_m \in B$ be a basis of $B$. $\text{Span}(v_1, \ldots, v_m) = \{\lambda_1 \land v_1 \lor \ldots \lor \lambda_m \land v_m \mid \lambda_1, \ldots, \lambda_m \in \mathbb{Z}_2\}$.

1 Throughout this paper we will use indistintely 0, F or false for the binary number 0 and 1, T or true for the binary number 1.

2 Bit strings can represent any complex data type. Consequently, our functional setting includes any functional computer program having fixed length input and output.
The dimension \( \dim(B) \) of the Boolean algebra is the minimum \( m \) such that 
\[ B = \text{span}(v_1, \ldots, v_m). \]

The method proposed in this paper consists of two combined strategies: reducing the dimensionality of the Boolean algebra where a target function lies and combining known operations in the algebra, using them as a basis to synthesize a program that emulates the target function.

Both strategies are implemented using two recursive higher order functions that we created and that we named the drill and the join.

In computer science higher-order functions are functions that take one or more functions as input and output a function. They correspond to linear mappings in mathematics.

### 4 The Drill function

We define the set \( F_m \) of functions \( f : \mathbb{Z}_2^p \times \mathbb{Z}_2^q \rightarrow \mathbb{Z}_2 \) containing Boolean functions belonging to a Boolean algebra of dimension \( m \), described in polynomial form:

\[
f(X, Y) = \bigoplus_{i=1}^{m} g_i(X) \land h_i(Y) \quad (2)
\]

where \( g_i : \mathbb{Z}_2^p \rightarrow \mathbb{Z}_2 \) and \( h_i : \mathbb{Z}_2^q \rightarrow \mathbb{Z}_2 \) are also Boolean functions. Note that equations 1 and 2 are equivalent and interchangeable. Equation 2 only splits its input space in two disjoint subsets: \( p + q = n \).

Considering a function \( f \in F_m \), a chosen \( X_0 \in \mathbb{Z}_2^p \) and a chosen \( Y_0 \in \mathbb{Z}_2^q \) such that \( f(X_0, Y_0) \neq 0 \), we define the drill higher-order function:

\[
\mathcal{F}_{X_0,Y_0}(X,Y) = \mathcal{F}(f(X,Y), X_0, Y_0) = f(X,Y) \oplus (f(X_0,Y) \land f(X,Y_0)) \quad (3)
\]

Note that the function \( \mathcal{F} \) outputs a new function and has as inputs the function \( f \) and instances of \( X \) and \( Y \), defining a position on \( f \) input space.

**Theorem:** If \( f \in F_m \) and \( f(X_0, Y_0) \neq 0 \), then \( \mathcal{F} \in F_r \) and \( r \leq m - 1 \).

**Proof:** Consider \( W = \text{span}(h_1, \ldots, h_m) \). Consequently \( \dim(W) \leq m \). The linear operator \( h \in W \rightarrow h(Y_0) \) is not the zero map because the hypothesis forbids \( h_i(Y_0) = 0 \) for all \( i = 1, \ldots, n \). Consequently, the vector subspace \( \overline{W} = \{ h \in W | h(Y_0) = 0 \} \) has \( \dim(\overline{W}) \leq m - 1 \). Notice that for all \( X \in \mathbb{Z}_2^p \) we have \( \overline{f}(X, \cdot) \in \overline{W} \). In fact:

\[
\overline{f}(X,Y_0) = f(X,Y_0) \oplus (f(X_0,Y_0) \land f(X,Y_0)) = 0 \quad (4)
\]

Let \( r = \dim(\overline{W}) \) and \( \overline{h}_i, i = 1, \ldots, r \) be a spanning set such that \( \overline{W} = \text{span}(\overline{h}_1, \ldots, \overline{h}_r) \). For all \( X \in \mathbb{Z}_2^p \), \( \overline{f}(X, \cdot) \) can be represented as a linear combination of the \( \overline{h}_i \), the coefficients depending on \( X \). In other words, there exist coefficients \( \overline{g}_i(X) \) such that:

\[
\overline{f}(X, \cdot) = \bigoplus_{i=1}^{r} \overline{g}_i(X) \land \overline{h}_i \quad (5)
\]
or written differently and remaining that \( r = \dim(\mathcal{W}) \) and consequently \( r \leq m - 1 \):

\[
\mathcal{J}(X, Y) = \bigoplus_{i=1}^{r} \pi_i(X) \land \bar{h}_i(Y) \quad \Box \tag{6}
\]

As an illustration let us consider a Boolean function \( f : \mathbb{Z}_2 \times \mathbb{Z}_2 \rightarrow \mathbb{Z}_2 \) whose behavior can be described by table 1. One possible representation for this function

<table>
<thead>
<tr>
<th>( x )</th>
<th>( y )</th>
<th>( f(x, y) )</th>
</tr>
</thead>
<tbody>
<tr>
<td>F</td>
<td>F</td>
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<td>T</td>
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</tbody>
</table>

would be \( f(x, y) = y \lor (-x \land -y) \). \( f \) is of dimension 2 (no shorter representation is possible). Note that this representation is given here for the illustration’s sake. The method does not require high-level definitions, only I/O examples.

Respecting the stated hypothesis we can pick \( x_0 = F \) and \( y_0 = F \) once \( f(F, F) = T \). The partial functions obtained will be: \( f(x_0, y) = y \lor (T \land -y) \) and \( f(x, y_0) = F \lor (-x \land T) \). Appplying the drill function we obtain:

\[
\mathcal{J}(x, y) = \mathcal{H}(f(x, y), x_0, y_0) = (y \lor (-x \land -y)) \oplus ((y \lor (T \land -y)) \land (F \lor (-x \land T))) = (x \land y) \tag{7}
\]

We can see that \( \mathcal{J}(x, y) \) is of dimension 1, confirming the stated theorem.

5 The Join function

Consider now the set \( F_m \) of Boolean functions \( f : \mathbb{Z}_2^n \rightarrow \mathbb{Z}_2 \) and \( v_1, \ldots, v_m \in F_m \) a basis. The functions in this set can be described in polynomial form as:

\[
f(X) = \bigoplus_{i=1}^{m} \lambda_i \land v_i(X) \tag{8}
\]

where \( \lambda_i \in \mathbb{Z}_2 \) are the coefficients.

Considering a function \( f \in F_m \), a chosen \( X_j \in \mathbb{Z}_2^n \) such that \( f(X_j) \neq 0 \) and a chosen function \( v_j \) belonging to the basis such that \( v_j(X_j) \neq 0 \), we define the join function:

\[
\mathcal{H}_{X_j, v_j} = \mathcal{H}(f(X), X_j, v_j) = f(X) \oplus v_j(X) \tag{9}
\]

**Theorem:** If \( f \in F_m \), \( f(X_j) \neq 0 \) and \( v_j(X_j) \neq 0 \), then \( \overline{f} = \mathcal{H}_{X_j, v_j} \in F_r \) and \( r \leq m - 1 \).
Proof: Consider $W = \text{span}(v_1, \ldots, v_m)$. Consequently $\dim(W) \leq m$. The linear operator $v \in W \rightarrow v(X_j)$ is not the zero map otherwise $v_j(X_j) = 0$. Consequently, the vector subspace $\overline{W} = \{ f \in W | f(X_j) = 0 \}$ has $\dim(\overline{W}) \leq m - 1$. We can see that $\overline{f} \in \overline{W}$. In fact:

$$\overline{f}(X_j) = f(X_j) \oplus v_j(X_j) = 0 \quad (10)$$

Let $r = \dim(\overline{W})$ and $\overline{v}_i, i = 1, \ldots, r$ be a spanning set such that $\overline{W} = \text{span}(\overline{v}_1, \ldots, \overline{v}_r)$. The function $\overline{f}$ can be represented as a linear combination of the $\overline{v}_i$.

$$\overline{f}(X) = \bigoplus_{i=1}^{r} \overline{v}_i \wedge v(X) \quad \square \quad (11)$$

We can use the same function $f : \mathbb{Z}_2 \times \mathbb{Z}_2 \rightarrow \mathbb{Z}_2$ described by table 1 to illustrate the behavior of the join higher-order function. $f$ belongs to a Boolean algebra of dimension 2 which can be defined, for instance, by the following spanning set: $v_1(x, y) = x, v_2(x, y) = y, v_3(x, y) = x \land y, v_4(x, y) = T$. Respecting the stated hypothesis we can pick $X_j = (T, T)$ and $v_j = v_1$ once $f(T, T) = T$ and $v_1(T, T) = T$. Applying the join function we obtain:

$$\overline{f}(x, y) = \overline{H}(f(X), X_j, v_j) = (y \lor (\neg x \land \neg y)) \oplus x = (x \land y) \quad (12)$$

We can see that $\overline{f}(x, y)$ is of dimension 1, confirming the stated theorem.

6 The Drill & Join program synthesis method

Drill and join are used to define a program synthesis method. Considering an active learning framework, the input function $f(X, Y)$ on $F$ and the input function $f(X)$ on $H$ represent an external unknown concept from which it is possible to obtain data by means of queries (input-output examples).

This unknown concept could be, for instance, some physical phenomenon that a machine with sensors and actuators can actively experiment, an algorithm to be translated in a hardware description language, a computer program that one would like to emulate or optimize or a decision process to be implemented on a computer for which one or more experts are able to answer required questions.

In order to understand the method it is important to notice two important properties of both higher-order functions:

- $F$ and $H$ can be applied recursively: if $f(X, Y) \in F_m$ then $\overline{f}_1(X, Y) = \overline{F}(f(X, Y), X_0, Y_0) \in F_{m-1}$ and $\overline{f}_2(X, Y) = \overline{H}(\overline{f}_1(X, Y), X_1, Y_1) \in F_{m-2}$.

Similarly, if $f(X) \in F_m$ then $\overline{f}_1(X) = \overline{H}(f(X), X_0, v_0) \in F_{m-1}$ and $\overline{f}_2(X) = \overline{H}(\overline{f}_1(X), X_1, v_1) \in F_{m-2}$. Each recursion generates a new function belonging to an algebra of a lower dimension.
The recursion ends when the higher-order functions become the zero map: 
\[ F(f(X, Y), X_i, Y_j) = 0 \Rightarrow f(X, Y) = (f(X, Y) \land f(X, Y)) \] and similarly, 
\[ H(f(X), X_i, v_i) = 0 \Rightarrow f(X) = v_i(X). \]

The first property enables us to apply the same higher order function recursively in order to gradually reduce the dimensionality of the initial problem. The second defines a stop condition. As a result, the output program is obtained.

\[ \forall X, Y \in \mathbb{Z}_2^n \times \mathbb{Z}_2^n : \overline{f}_m(X, Y) = F(\overline{f}_{m-1}(X, Y), X_m, Y_m) = 0 \Rightarrow \overline{f}_m(X, Y) = (\overline{f}_{m}(X_m, Y) \land \overline{f}_m(X, Y_m)) \]

Replacing \( \overline{f}_m(X, Y) = F(\overline{f}_{m-1}(X, Y), X_m, Y_m-1) \) gives us:
\[ \overline{f}_{m-1}(X, Y) \oplus (\overline{f}_{m-1}(X_m-1, Y) \land \overline{f}_{m-1}(X, Y_m-1)) = (\overline{f}_{m}(X_m, Y) \land \overline{f}_m(X, Y_m)) \] (13)
and consequently:
\[ \overline{f}_{m-1}(X, Y) = (\overline{f}_{m-1}(X_m-1, Y) \land \overline{f}_{m-1}((X, Y_m-1)) \oplus (\overline{f}_{m}(X_m, Y) \land \overline{f}_m(X, Y_m)) \] (14)

Tracking the recursion back to the beginning we obtain:
\[ f(X, Y) = \bigoplus_{i=1}^{m} \overline{f}_i(X_i, Y) \land \overline{f}_i(X, Y) \] (15)

Equation 15 tells us that the original target function \( f \) can be recreated using the partial functions \( \overline{f} \) obtained using the drill function. The partial functions \( \overline{f} \) are simpler problems, defined on subspaces of the original target function.

Similarly:
\[ \forall X \in \mathbb{Z}_2^n : \overline{H}(\overline{f}_{m}(X), X_m, v_m) = 0 \Rightarrow \overline{f}_m(X) = v_m(X) \] (16)

Replacing \( \overline{f}_m(X) = H(\overline{f}_{m-1}(X), X_m-1, v_{m-1}) \) gives us:
\[ \overline{f}_{m-1}(X) \oplus v_{m-1}(X) = v_m(X) \] (17)
and:
\[ \overline{f}_{m-1}(X) = v_{m-1}(X) \oplus v_m(X) \] (18)

Tracking the recursion back to the beginning we obtain:
\[ f(X) = \bigoplus_{i=1}^{m} v_i(X) \] (19)

Equation 19 tells us that the original target function \( f \) can be recreated using the partial functions \( \overline{f} \) obtained using the join function and the basis \( v \).

Note that if the drill initial condition cannot be established, i.e., no \( (X_0, Y_0) : f(X_0, Y_0) \neq 0 \) can be found, the target function is necessarily \( f(X, Y) = F \). On the same way if no \( X_0 : f(X_0) \neq 0 \) can be found to initiate join the target function is the zero map \( f(X) = F \). If it exists a \( X_j : f(X_j) \neq 0 \) but no \( v_j : v_j(X_j) \neq 0 \) the basis was not chosen appropriately.
The IF higher order function defines a double recursion. For each step of the dimensionality reduction recursion two subspace synthesis problems are defined: $\mathcal{F}(X_1, Y)$ and $\mathcal{F}_f(X, Y)$. Each of these problems can be treated as a new target function in a Boolean algebra of reduced arity once part of the arguments is fixed. They can be recursively solved using IF or IH again.

The combination of IF and IH defines a very powerful inductive method. The IF higher-order function alone requires the solution of an exponential number of subspace synthesis problems in order to inductively synthesize a target function. The IH higher-order function requires a basis of the Boolean algebra. Considering the $2^n$ cardinality of a basis, it can be impractical to require its prior existence in large arity algebras. Nevertheless, both functions combined can drastically reduce the number of queries and the prior bases definition.

The method, detailed on the next sections, uses the following strategies:
- Bases are predefined for low arity input spaces, enabling the join function.
- Synthesis on large arity input spaces begin using the IF function. 
- Previously synthesized programs on a subspace can be memorized in order to compose a basis on that subspace.
- At each new synthesis, if a basis exists use IH, otherwise use IF.

To illustrate the functioning of the whole method let us use again the same function $f : \mathbb{Z}_2 \times \mathbb{Z}_2 \rightarrow \mathbb{Z}_2$ described by tables 1 and 2. We can apply the drill function one first time $\mathcal{F}_1(X, Y) = \mathcal{F}(f(X, Y), X_0, Y_0)$ with $x_0 = F$, $y_0 = F$, defining two subspace problems: $f(x_0, y)$ and $f(x, y_0)$. Both problems can be solved using the join function and the basis $v_0(x) = x$, $v_1(x) = T$. $\mathcal{F}_1$ is not the zero map, requiring a second recursion of drill $\mathcal{F}_2 = \mathcal{F}(\mathcal{F}_1(X, Y), X_1, Y_1)$ with $x_1 = T$, $y_1 = T$. Two new subspace problems are defined: $\mathcal{F}_1(x_1, y)$ and $\mathcal{F}_1(x, y_1)$ and they can be solved using join and the same basis again. $\mathcal{F}_2(X, Y) = F$ is finally the zero map, stopping the recursion. Table 2 shows the target function and the transformation steps performed using the drill function. To illustrate

<table>
<thead>
<tr>
<th>x</th>
<th>y</th>
<th>f(x, y)</th>
<th>f(x_0, y)</th>
<th>f(x, y_0)</th>
<th>f_1(x, y)</th>
<th>f_1(x_1, y)</th>
<th>f_1(x, y_1)</th>
<th>f_2(x, y)</th>
</tr>
</thead>
<tbody>
<tr>
<td>F</td>
<td>F</td>
<td>T</td>
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<td>T</td>
<td>F</td>
</tr>
</tbody>
</table>

Table 2. Truth table for the drill steps

the use of the join function let us consider the reconstruction of $f(X) = f(x, y_0)$ detailed in table 3. One first application $\mathcal{H}(X) = \mathcal{H}(f(X), X_0, v_0)$ with $x_0 = F$ and $v_0(x) = x$ will not result in the zero map, as shown on the fifth column of table 3, requiring a recursive call $\mathcal{H}_2(X) = \mathcal{H}(\mathcal{H}(X), X_1, v_1)$ with $x_1 = T$ and $v_1(x) = T$ which will result the zero map.
Using equation 19 we can find a representation for the partial target function:
f(X) = f(x, y_0) = v_0(x) \oplus v_1(x) = x \oplus T. The same process can be used to find representations for all one-dimensional problems: f(x_0, y) = T, f_1(x_1, y) = y and f_1(x, y_1) = x.

Having solved the partial problems we can use equation 15 to build the full target function:
\[ f(x, y) = (f(x_0, y) \land f(x, y_0)) \oplus (f_1(x_1, y) \land f_1(x, y_1)) = (T \land (x \oplus T)) \oplus (y \land x) \text{ which is equivalent to our initial representation.} \]

Each higher order function underlies a query protocol. At each recursion of If one or more queries for data are made in order to find \(X_i, Y_i \in \mathbb{Z}_2 \times \mathbb{Z}_2\): \(f_i(X_i, Y_i) \neq 0\). At each recursion of If a position \(X_i \in \mathbb{Z}_2 : f(X_i) \neq 0\) and a function \(v_i : v_i(X_i) \neq 0\) from the basis are chose. The queries require data from the target function and must be correctly answered by some kind of oracle, expert, database or system. Wrong answers make the algorithms diverge. Data is also necessary to verify the recursion stop condition. A full test of \(f_i(X, Y) = 0\) or \(f_1(X) = 0\), scanning the whole input space, can generate proof that the induced result exactly meets the target function. Partial tests can be enough to define candidate solutions subject to further inspection.

7 The main algorithms of the Drill&Join method

To explain how the drill and join higher-order functions can be used as program-forming functionals we propose the following two algorithms.

Drill takes as initial inputs the target function: \(f_n\) and the dimension of its input space: inputs. Deeper inside the recursion \(f_n\) corresponds to \(f_{m}(X, Y)\), inputs defines the dimension of its subspace and initial indicates its first free dimension. Drill returns a synthesized functional program that emulates \(f_n\).

1: procedure Drill\(\text{(}f_n, \text{inputs,optional: initial = 0})\)
2: if have a basis for this subspace then
3:    return Join\(\text{(}f_n, \text{inputs,initial})\);
4: end if
5: pos \leftarrow \text{find a position inside this subspace where } f_n \text{ is not null};
6: if pos = null then
7:    return FALSE; \quad \triangleright \text{Stop condition.}
8: end if
9: fa(args) = fn(concatenate(args, secondhalf(pos))); \quad \triangleright f(X, Y_0)
10: fb(args) = fn(concatenate(firsthalf(pos), args)); \quad \triangleright f(X_0, Y)
11: fc(args) = fn(args) \oplus (fa(firsthalf(args)) \land fb(secondhalf(args)));
12: \quad \triangleright f_{m+1}(X, Y)
13: pa = Drill\(\text{(}fa, \text{inputs}/2, \text{initial})\); \quad \triangleright \text{Recursive call to synthesize } f(X, Y_0)
14: \( pb = \text{Drill}(fb, \text{inputs}/2, \text{initial} + \text{inputs}/2); \)  \( \triangleright \) Recurs. call for \( f(X_0, Y) \)
15: \( pc = \text{Drill}(fc, \text{inputs}, \text{initial}); \)  \( \triangleright \) Recurs. call for \( f_{m+1}(X, Y) \)
16: \( \text{return } \left( \text{'( 'pa 'AND' pb ')'} \text{ XOR } pc \text{' )'} \right); \)  \( \triangleright \) Returns the program
17: \textbf{end procedure}

Note that \( fa, fb \) and \( fc \) are functions based on \( fn \) and \( pa, pb \) and \( pc \) are programs obtained by recursively calling \( \text{drill} \), \text{firsthalf} \ and \text{seconhalf} split a vector in two halves.

\( \text{Join} \) takes as input a function \( fn \), belonging to a Boolean algebra. The algorithm requires a basis for this Boolean algebra, materialized as an array of functions: \text{basis}[] \ and an array of programs emulating the basis: \text{basisp}[]. The algorithm creates a program to emulate \( fn \) combining the program basis.

1: \textbf{procedure} \text{Join}(fn, \text{optional: initial = 0})
2: \( \text{pos} \leftarrow \text{find a position inside this subspace where fn is not null}; \)
3: \textbf{if} \( \text{pos} = \text{null} \) \textbf{then}
4: \( \text{return FALSE}; \)  \( \triangleright \) Stop condition.
5: \textbf{end if}
6: \( v \leftarrow \text{find a function v from the basis such that v(pos) is not 0}; \)
7: \( vp \leftarrow \text{get the program that emulates v}; \)
8: \( fa(args) = fn(args) \oplus v(args); \)
9: \( pa = \text{Join}(fa, \text{initial}); \)  \( \triangleright \) Recursive call for \( f_{m+1}(X) \)
10: \( \text{return } \left( \text{'( 'pa 'XOR' vp(initial) ')'} \right); \)  \( \triangleright \) Returns the program.
11: \textbf{end procedure}

8 \textbf{A Common Lisp version of the Drill\&Join method}

Common Lisp is a natural choice of programming language to implement the \textit{drill\&join} method. Lisp functions can take other functions as arguments to build and return new functions. Lisp lists can be interpreted and executed as programs.

The following code implement the simplest version of the method. It synthesizes a Lisp program that emulates a function \( fn \) by just querying it. The function \( fn \) must accept bit strings as inputs and must return one bit as the answer. In this simple illustration \( fn \) is another Lisp function but in real use it would be an external source of data queried throughout an adequate experimental protocol.

\( \text{Drill} \) takes the unknown function \( fn \) and its number of binary arguments \( nargs \) as input. It returns a list composed of logical operators, logical symbols \text{nil} \ and \text{true} \ and references to an input list of arguments. The output list can be executed as a Lisp program that emulates \( fn \).

\begin{verbatim}
(defun drill (fn nargs &optional (ipos 0) (slice 0))
  (let ((base (list nil #'(lambda(x) (first x)) #'(lambda(x) t) ))
        (basep (list nil #'(lambda(x) (list 'nth x 'args)) #'(lambda(x) t))))
    (if (= nargs 1)
        (join fn base basep ipos)
        (let ((pos (findpos fn nargs slice)))
          (if (null pos)
              nil
              (labels ((fa (args) (funcall fn (append args (cdr pos))))
                        (fb (args) (funcall fn (append (list (car pos)) args)))
                        (fc (args) (xor (funcall fn args)))))
            (join fn base basep ipos) (pos (findpos fn nargs slice)))
          (labels ((fa (args) (funcall fn (append args (cdr pos))))
                        (fb (args) (funcall fn (append (list (car pos)) args)))
                        (fc (args) (xor (funcall fn args))))
            (join fn base basep ipos) (pos (findpos fn nargs slice)))
          (labels ((fa (args) (funcall fn (append args (cdr pos))))
                        (fb (args) (funcall fn (append (list (car pos)) args)))
                        (fc (args) (xor (funcall fn args))))
            (join fn base basep ipos) (pos (findpos fn nargs slice))))
  (labels ((fa (args) (funcall fn (append args (cdr pos))))
                        (fb (args) (funcall fn (append (list (car pos)) args)))
                        (fc (args) (xor (funcall fn args))))
            (join fn base basep ipos) (pos (findpos fn nargs slice))))
)
\end{verbatim}
(and (fa (list (car args))) (fb (cdr args))))

(let ((sa (drill #'(lambda(args) (fa args)) 1 ipos))
  (sb (drill #'(lambda(args) (fb args)) (1- nargs) (1+ ipos)))
  (sc (drill #'(lambda(args) (fc args)) nargs ipos (1+ slice))))
  (if (and (atom sa) sa)
    (setq r1 sb)
    (setq r1 (list 'and sa sb)))
  (if (null sc)
    r1
    (list 'xor r1 sc))))))))

In this implementation the split of the input space is done by choosing the first argument to be $X$ and the rest to be $Y$. Drill will call Join when the recursion is down to just one input function ($nargs=1$). A basis for one input bit functions $\{f(x) = x, f(x) = T\}$ is defined directly inside drill as a functions list base and a programs list basep which are passed as arguments to the join function which is implemented as follows:

(defun join (fn base basep &optional (ipos 0))
  (let ((pos (findpos fn 1)))
    (if (null pos)
      nil
      (let ((fb (findbase base basep pos)))
        (labels ((fa (args) (xor (funcall fn args) (funcall (nth 0 fb) args))))
          (let ((r (join #'(lambda(args) (fa args)) base basep ipos)))
            (return-from join (list 'xor (funcall (nth 1 fb) ipos) r))))))))

The findpos function is used to test if fn is the zero map performing a full search on the fn input space. It stops when a non-zero answer is found and returns its position. The full search means that no inductive bias was used.

(defun findpos (fn nargs)
  (loop for i from 0 to (1- (expt 2 nargs)) do
    (let ((l (make-list nargs)) (j i) (k 0))
      (loop do (if (= (mod j 2) 1) (setf (nth k l) t))
        (incf k) (setq j (floor j 2))
      while (> j 0))
    (if (funcall fn l) (return-from findpos l)))))

The findbase function is used inside join to find a function from the basis respecting the constraint $v_j(X_j) \neq 0$.

(defun findbase (base basep pos)
  (loop for i from 1 to (1- (list-length base)) do
    (if (funcall (nth i base)pos)
      (let ((ba (nth i base)) (bp (nth i basep)))
        (setq fb (list ba bp))
        (delete (nth 0 fb) base) (delete (nth 1 fb) basep)
        (return-from findbase fb))))))

The method queries the target function (funcall fn) only inside findpos, in order to find a non-zero position.

Using the Lisp code provided above it is possible to check the consistency of the method. Applied to our illustration described by table 1 the generated program would be:

(XOR (AND (XOR T (NTH 0 ARGS)) T) (AND (NTH 0 ARGS) (NTH 1 ARGS)))
As a more advanced illustration, let us consider the "unknown" target function to be the Fibonacci sequence. To avoid bulky outputs we will limit the illustration to have as input an unsigned integer between 0 and 63. Consequently, the input of the \textit{fn} function can be a six bit long bit string. The range of the output (between 1 and 6557470319842) requires a 64 bits unsigned integer. The target function \texttt{fibonacci(n)} computes a long integer corresponding to the \textit{n}–\textit{th} position of the Fibonacci sequence. To translate integers to lists of \{'\texttt{NIL}', \texttt{T}'\} handled by the \textit{drill} and \textit{join} lisp code we use their binary representation. The translation is done by the routines \texttt{longint2bitlist}, \texttt{bitlist2longint}, \texttt{6bitlist2int} and \texttt{6int2bitlist}, not included in this paper, called from the function \texttt{myfibonacci(n)}.

\begin{verbatim}
(defun myfibonacci(n) (let ((r (6bitlist2int n))) (longint2bitlist (fibonacci r))))
\end{verbatim}

In order to synthesize a program able to emulate the whole target function we need to call \textit{Drill} for each output bit and generate a list of boolean expressions.

\begin{verbatim}
(defun synthesis (fn nargs nouts filename)
  (with-open-file (outfile filename :direction :output)
    (let ((l (make-list nouts)))
      (loop for i from 0 to (1- nouts) do
        (labels ((fa (args) (nth i (funcall fn args))))
          (let ((x (drill #'(lambda(args) (fa args)) nargs))) (setf (nth i l) x)))))
    (print l outfile)))
\end{verbatim}

To run the generated program we need to compute each output bit and translate the bit string into an integer:

\begin{verbatim}
(defun runprogram(filename v)
  (with-open-file (infile filename)
    (setq s (read infile)) (setq args (6int2bitlist v))
    (let ((r (make-list (list-length s))))
      (loop for i from 0 to (1- (list-length s)) do
        (setf (nth i r) (eval (nth i s))))
      (print (bitlist2longint r))))
\end{verbatim}

A full check on the whole target function input space shows that the synthesized program exactly emulates the target function. To illustrate how the generated programs looks like we show below the expression that computes the first output bit of the Fibonacci sequence:

\begin{verbatim}
(XOR (AND (XOR (AND (XOR (AND (NTH 0 ARGS) (XOR T (NTH 1 ARGS))) (XOR T (NTH 0 ARGS))) (NTH 1 ARGS))) T) (AND (XOR (AND (XOR T (NTH 0 ARGS)) (NTH 1 ARGS))) (XOR T (NTH 0 ARGS))) (NTH 2 ARGS)))) (AND (XOR (AND (XOR (AND (NTH 1 ARGS) (XOR T (NTH 0 ARGS))) T) (AND (XOR (AND (XOR T (NTH 0 ARGS)) (NTH 1 ARGS))) (XOR T (NTH 2 ARGS))) (AND (XOR (AND (NTH 0 ARGS) (XOR T (NTH 0 ARGS))) (XOR T (NTH 1 ARGS))) (AND (XOR T (NTH 0 ARGS)) (NTH 1 ARGS))) (NTH 2 ARGS))) (AND (XOR (AND (XOR T (NTH 3 ARGS)) (XOR T (NTH 4 ARGS))) (XOR T (NTH 5 ARGS))) (AND (XOR (AND (NTH 3 ARGS) (XOR T T)) (AND (XOR T (XTH 4 ARGS))) (XOR T (NTH 5 ARGS))) (AND (XOR (AND (XOR (AND (NTH 3 ARGS) (XOR T T)) (AND (XOR T (NTH 4 ARGS)) (XOR T (NTH 5 ARGS))) (XOR T (NTH 2 ARGS))) (AND (XOR (AND (NTH 0 ARGS) (XOR T T)) (AND (XOR T (NTH 1 ARGS))) (AND (XOR T (NTH 0 ARGS)) (NTH 1 ARGS))) (NTH 2 ARGS))) (AND (XOR (AND (XOR T (NTH 3 ARGS)) (XOR T (NTH 4 ARGS))) (XOR T (NTH 5 ARGS))))) (AND (XOR (AND (XOR T (NTH 3 ARGS)) (XOR T (NTH 4 ARGS))) (XOR T (NTH 5 ARGS)))))))
\end{verbatim}

The generated program is basically a single Boolean expression that explicitly references an input list called \texttt{args}. The full program consists of a list of those Boolean expressions.
9 Discussion

The implementation presented on section 8 has didactic purposes only. A number of enhancements are possible. Bases can be dynamically built on subspaces by memorizing programs previously created avoiding drilling down to smaller subspaces and reducing the number of queries to the target function.

An interesting aspect of the generated programs is the fact that the computation of each output bit is completely independent of the others. There is a specific program for each output bit, enabling parallel processing.

Empirical comparisons between the proposed method and existing ones are difficult because of the conceptual differences between them. Existing benchmark frameworks, as CHStone [17] and SyGuS [13] tend to be specific to a certain synthesis approach. CHStone was conceived for C-based high-level synthesis and SyGuS for syntax-guided hybrid synthesis (inductive and deductive). Conceptual comparisons can be done but without objective results. The proposed method can handle target functions that others probably cannot, but because it works at the bit level of inputs and output, the number of examples required for learning and testing tend to be larger than in other methods. Most existing inductive methods use static example databases while our method is based on active learning and requires an experimental protocol in order to query the target concept during the learning process. Our method requires a predefined input space, with fixed length, and does not handle dynamic input lists like in Lisp programs generated by most variants and extensions of Summers [9] analytical approach. But on the other side, the simplicity of the code generated by our method, based on only two logical operators, enables its compilation in almost any conventional programming language, even on hardware description languages. The declarative nature of the generated programs brings predictability at runtime in terms of execution time and use of machine resources. The example proposed in section 8 showed how to synthesize a declarative program to emulate the Fibonacci sequence. The generated program requires, for any six bits input, exactly 1804 low-level, bitwise logic operations to compute the corresponding Fibonacci number. An equivalent imperative program will require local variables, loop or recursion controls and a number of variable assignments and arithmetic operations proportional to the input value. The method does not require any prior knowledge about the target function, as a background theory, types of variables or operation on types, like in hybrid methods. Nevertheless, the goal of this paper is not to prove that we created a better method but to present a new concept on inductive program synthesis. Future work will be necessary in order to assess the advantages and disadvantages of the proposed method in each possible field of application when other methods are also available.

For practical use of our method it is important to be able to estimate the effort required to synthesize a program. In typical active learning, there is usually a cost element associated with every query. This cost depends on the characteristics of the target concept and the associated experimental protocol used to query it. The synthesis effort will depend on the number of queries to be required and the cost of each query. If the implementation of the method is based on a full
verification of the zero map, like in our illustration presented on section 8, a full scan of the target function input space will be necessary and the number of queries will depend only on the size of the input space in bits. As a future work an inductive bias can be proposed in order to avoid the full scan and then reduce the number of queries.

10 Applications

Applications of this method include all sorts of systems based on bitwise operations, given that the learning problem can be described in functional form. We successfully applied the method to target functions handling different complex datatypes, as floating-point numbers and texts.

Any functional computer program can be emulated using our approach. There can be a number of reasons to perform reverse engineering of existing computer programs: the lack of a source code, the need to translate an executable program to run on a different computational platform, the intent to optimize an inefficient implementation.

The method can also be used to translate algorithms from a high-level language directly into combinatorial circuit design expressed in an hardware description language. The code generated by the proposed method can be easily mapped into a netlist (sequence of circuit gates).

Machine learning is another field of application. Machines having sensors and actuators, like robots, can acquire direct models using the method. The machine can actively experiment a physical phenomena and synthesize a program to predict the result of possible actions.

11 Conclusion

We have presented Drill&Join, a generic method that actively interacts with an external concept (system, function, oracle or database) via I/O examples, and synthesizes programs. Generated programs are based on Boolean expressions. The method is not restricted to any specific form of functional learning problem or target function and does not require any background knowledge to be applied. The only requirement is that the external source of data is consistent. Our work presents a number of interesting questions for future consideration. The combination of the drill and the join higher-order functions and the dynamic construction of bases on subspaces via memorization of generated program can drastically reduce the number of recursive calls and queries. Further investigation is necessary on how to explore these dynamic bases on large input spaces. The stop condition of the algorithms, based on a full verification of the zero map, requires a complete scan of the learning input space. Partial verifications can compromise the convergence of the algorithms. Investigations on inductive biases adequate to the method are necessary.
References

Automatic Synthesis of Combiners in the MapReduce Framework
An Approach with Right Inverse

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Abstract. We give an algorithm for synthesizing combiners, which partially aggregate data inside mappers before sending to reducers, to improve communication cost of MapReduce programs. Our algorithm, given a pair of a mapper and a reducer, automatically generates a right inverse of a reducer and composes it with the reducer. Then the composed function can be used both as a reducer and a combiner. We implement our algorithm and conduct experiment in Amazon Elastic MapReduce. The result is promising: Generated combiners reduce significant amount of data.

Keywords: Program transformation, MapReduce, Right inverse, Data processing

1 Introduction

MapReduce [3] is a framework that makes distributed computation accessible. In the MapReduce framework, a programmer writes a program so that it consists of two steps: MAP and REDUCE. Then, the framework automatically conducts distributed data processing.

Figure 1 shows how a MapReduce computation proceeds. At the MAP step, the input data are distributed to mappers to be converted to a list of pairs of a key and a value. Then, pairs with an identical key are gathered to a single reducer where the pairs are reduced to a single output. For example, a MapReduce program that counts the frequency of each word in text files works as follows. A mapper, for each occurrence of a word $w$, generates a key–value pair $(w, 1)$ and send it to a reducer associated with $w$. Each reducer counts the number of received key–value pairs and outputs the result.

Mapper and reducer processes are automatically distributed over mapper nodes and reducer nodes, respectively, and the computation in different mapper nodes and reducer nodes proceeds in parallel. Thereby, a programmer, if her program can be expressed in these two steps, does not need to be bothered by difficulties in distributed programming such as synchronization and communication.

One optimization technique for MapReduce computation is reduction of the communication cost between mapper and reducer nodes. In general, a reducer
works at a node distant from mappers. Hence, reducing the amount of data communicated to reducers is important for performance tuning.

A method to address this problem is the use of combiners. A combiner works in each mapper node and conducts the computation that is to be conducted in a reducer before data are transferred. By conducting the reducer-like computation in advance, a combiner often decreases the amount of transferred data. However, as we discuss later, it is not always the case that a reducer can be used as a combiner without change.

We propose a method for automatically generating a combiner from a mapper and a reducer. Concretely, we propose the following program transformation:

**Input:** A mapper $m$ and a reducer $r$.

**Output:** A mapper $m'$ and a reducer $r'$, which also serves as a combiner.

The mapper $m'$, combiner $r'$, and reducer $r'$, as a whole Map-Reduce program, implement the same computation as $m$ and $r$.

Roughly speaking, our algorithm automatically derives a weak right inverse\footnote{A weak right inverse of a function $f : A \to B$ is a function $g : B \to A$ such that $f \circ g \circ f = f$.} (we often omit “weak” in what follows) of a given reducer and composes it with the reducer. Then, the composed function can be used both as a reducer and a combiner. In theory, the only condition for this transformation to work is that $r$ is a list homomorphism \cite{harper1990}. In our implementation, $r$ is written in a DSL \cite{harper1990} so that a right inverse is systematically derived.
Note that our theoretical framework assumes that the reducer function is not necessarily commutative or, in other words, can be sensitive to the order of data. It means that it can deal with a wider class of problems, including Maximum Prefix Sum (MSP). We have implemented our system so that order-sensitive computation can be conducted.

Our combiner-synthesis algorithm is useful in the following aspects.

- The generated combiner is guaranteed to be correct by construction. This is indeed important since the behavior of a combiner is highly non-deterministic; see Section 2 for detail.
- Our algorithm prevents code duplication between a reducer and a combiner; it is often the case that a combiner often shares much code with a reducer because the computation conducted by a combiner is essentially the same as that by a reducer.
- Synthesized combiners work well. Our experiment shows that generated combiners reduce a significant amount of communicated data.

The rest of the paper is structured as follows. Section 2 overviews the MapReduce framework. Section 3 presents our combiner-synthesis algorithm. Section 4 gives a proof of correctness of our algorithm, which is followed by description of the current implementation (Section 5) and by the result of the experiment (Section 6). Section 7 discusses related work and Section 8 concludes.

2 The MapReduce framework

This section describes the MapReduce framework and illustrates our combiner-synthesis method by examples. We first present an example of a MapReduce program. We use a Haskell-like syntax to present a program in the rest of this paper.

Example 1. We present a MapReduce program that calculates the sum of integers. The function mapper takes a string as input and parses it to obtain a key and an integer; the functions getKeyFromString and getValueFromString parse the input string.

mapper :: string -> (key, Integer)
mapper str = (getKeyFromString str, getValueFromString str)

Here, (key, Integer) is the type of pairs of a key and an integer. The input data are distributed to several mapper nodes. The pairs generated by mappers within a single mapper node are sorted and aggregated by the key part and sent

2 One may wonder if MapReduce assumes that the order of input data is irrelevant. It is true, but recent empirical study by Xiao et al. [12] has shown that such order-sensitive computation is not uncommon in MapReduce code base. As discussed in Section 7, Liu et al. [7], who propose a transformation from sequential programs to MapReduce programs, have also taken order-sensitive computation into account, although they do not consider combiners.
to reducers; this phase is conducted by shufflers that are provided by MapReduce implementations. For example, suppose that strings “A 3”, “A 5”, “B 6”, “A 7”, and “B 8” are fed to mappers. The mappers convert these strings to key-value pairs (“A", 3), (“A", 5), (“B", 6), (“A", 7), and (“B", 8). Then, a shuffler sorts and aggregates the pairs by each key (i.e., (“A", [3,5,7]) for key “A"; (“B", [6,8]) for key “B") and sends them to reducers.

The function reducer, given a key and a list of integers, calculates the sum of the integers,

\[
\text{reducer} :: (\text{key}, [\text{Integer}]) \rightarrow (\text{key}, \text{Integer})
\]

\[
\text{reducer} (k, vs) = (k, \text{sum } vs)
\]

Here, [Integer] is the type of lists of integers. The function \text{sum} receives a list of integers and calculates the sum of them. For the previous data, a reducer associated with the key “A” receives (“A", [3,5,7]) from the previous mapper node and data of the form (“A", [...]) from other mapper nodes. The reducer calculates the sum of all the integers and outputs it.

As mentioned in Section 1, we can enhance the efficiency of the program above by using combiners to reduce the amount of data transferred from mappers to reducers. The mapper, during its execution, calls this combiner with the generated key-value pairs to curtail the number of pairs incrementally. A combiner should not affect the result of the whole computation.

Example 2. A combiner for the program in Example 1 could be implemented as follows.

\[
\text{combiner} :: (\text{key}, [\text{Integer}]) \rightarrow (\text{key}, \text{Integer})
\]

\[
\text{combiner} (k, vs) = (k, \text{sum } vs)
\]

Suppose the program above is invoked with the same input as that in Example 1. One possible scenario of the combiner execution is as follows: (1) The pairs (“A", 3) and (“A", 5) generated by the mappers are passed to the combiner that returns (“A", 8), (2) The pairs (“B", 6) and (“B", 8) are passed to the combiner generating (“B", 14), and (3) The pairs (“A", 8), (“A", 7) are passed again to the combiner to generate (“A", 15). The resulting pairs (“A", 15) and (“B", 14) are aggregated by shuffler to (“A", [15]) and (“B", [14]) and then sent to reducers. The size of data transferred from mappers to reducers is reduced: from (“A", [3,5,7]) and (“B", [6,8]) to (“A", [15]) and (“B", [14]).

Another scenario with the input data above would be as follows: (1) assembling

3 The order of key-value pairs is not preserved by a shuffler in general. Hence, the computation conducted by a reducer would be best described by a function that takes a set of values rather than a list of values. We, however, see a reducer as a function that takes a list so that our framework allows order-sensitive computation. See Section 5 for more discussion.

4 Combiners are invoked when a buffer inside a mapper node is full. Thus, the timing of combiner invocations depends on, for example, the size of the buffer, input data, and scheduling of mapper execution.
("A", 3) and ("A", 5) to produce ("A", 8), (2) assembling ("B", 6) and ("B", 8) to produce ("B", 14), and (3) ("A", [8,7]) and ("B", [14]) are sent to reducers. A programmer has to take non-determinism into account in writing a combiner so that the final result of computation is the same independently of how combiners are invoked.

Although the reducer can be used as a combiner without any change in Example 2, it is not always the case.

**Example 3.** Suppose that, for each key, we are to calculate the average of integers that are associated with the key. A mapper, a reducer, and a combiner can be implemented as follows.

```haskell
mapper :: str -> (key, (Integer, Integer))
mapper str = (getKeyFromString str, (1, getValueFromString str))

combiner :: (key, [(Integer, Integer)]) -> (key, (Integer, Integer))
combiner (k, vs) = (k, (sum (map fst vs), sum (map snd vs)))

reducer :: (key, [(Integer, Integer)]) -> (key, Integer)
reducer (k, vs) = (k, total / freq)
  where (freq, total) = (sum (map fst vs), sum (map snd vs))
```

The key–value pairs communicated among these three functions are of type (key, (Integer, Integer)). The value part consists of a pair of the number of integers that are summed up (we call the number frequency) and the sum of the integers. The function `map` takes a function `f` and a list `l`, and returns a list that is obtained by applying `f` to each element of the list `l`.

The mapper, reading a key `k` and an integer `n`, returns `(k, (1, n))` because the sum `n` is obtained from one integer. The frequency parts of the received key–value pairs are summed up by the combiner and the reducer; the sum parts are also summed up by the combiner and the reducer. The reducer divides the sum part by the frequency part to obtain the average of all the integers associated to a key.

**Remark 1.** Although in MapReduce the order of data is not preserved basically, adding indices to data enables list-operations; the order of data can be restored. The details are explained in Section 5.

### 3 Combiner synthesis

This section presents how to synthesize a combiner. As mentioned in the introduction, our algorithm takes a mapper and a reducer and synthesizes a new mapper and a new reducer, which can be also used as a combiner.

Recall that a required property for a combiner is that, regardless of when and how many times it is executed, the final outcome does not change. To guarantee this property, we take advantage of the following property of right inverse: If `g`
is a right inverse of \( f \), then \( f(g(f(g(f(v)))) \) is equal to \( f(v) \)—in other words, \( f \circ g \) is idempotent for the image of \( f \). Roughly speaking, we take a reducer as \( f \) and use \( f \circ g \) as a new reducer/combiner. That is,

- the synthesized mapper applies the original mapper and then the original reducer to an input; and
- the synthesized reducer (combiner) applies the right inverse of the original reducer to each element of an input list and aggregates them using the original reducer.

Then, repeated applications of combiners will cancel. Now, we elaborate our synthesis.

Let \( m \) and \( h \) be the input mapper and reducer, respectively. We assume that their types are

\[
\begin{align*}
m &: \text{data} \to a \\
r &: [a] \to b
\end{align*}
\]

where \( \text{data} \) is the type of the input to the mapper. (For simplicity, we assume all the keys are identical in this section and so omit keys.)

First we generate a right inverse of \( r \); we write it \( \text{inv } r \), and the type of \( \text{inv } r \) is \( b \to [a] \). Several methods, including Morita et al. [9], are readily applicable for generating a right inverse from a given function.

Then, the output of our algorithm (i.e., a new mapper, a new combiner, and a new reducer) is defined as follows.

\[
\begin{align*}
\text{mapper} &: \text{data} \to b \\
\text{mapper } v &= r \ [m \ v] \\
\text{combiner} &: [b] \to b \\
\text{combiner } vs &= r \ \text{concat} \ (\text{map} \ (\text{inv } r) \ vs) \\
\text{reducer} &: [b] \to b \\
\text{reducer } vs &= r \ \text{concat} \ (\text{map} \ (\text{inv } r) \ vs)
\end{align*}
\]

The definitions of the auxiliary functions \( \text{map} \) and \( \text{concat} \) are standard; see Appendix for their definitions.

**Example 4.** We explain this method in the simple case: sum. When \( r = \text{sum} \), one option of \( \text{inv } \text{sum} \) is the function to return the singleton list consisting of an argument—that is, \( (\text{inv } \text{sum}) \ x = [x] \). It is easy to check this is indeed a right inverse. Then, the synthesized functions are:

\[
\begin{align*}
\text{mapper } v &= \text{sum} \ [m \ v] = m \ v \quad \text{-- we expect } m \text{ to return an integer} \\
\text{combiner } vs &= \\
\text{sum} \ (\text{concat} \ (\text{map} \ (\text{inv } \text{sum}) \ vs)) \quad \text{-- equal to } \text{sum } vs \\
\text{reducer } vs &= \text{sum } vs
\end{align*}
\]
However, our method cannot be applied to every reducer. Consider the following reducer, which is a more straightforward version of the average-calculating program in Example

```haskell
reducer :: [Integer] -> Integer
reducer vs = (sum vs) / (len vs)
```

where\(\text{len}\) is the function that returns the length of the given list. Then, a right inverse of this reducer is the same as \(\text{inv \ sum}\) above. The synthesized combiner/reducer (we call \(\text{new\_red}\)), however, does not work as expected. Indeed, \(\text{new\_red}[2,4,9] = 5\), but \(\text{new\_red}([\text{new\_red}[2,4], \text{new\_red}[9]]) = 6\). To make our synthesis to work, we can use the following mapper and reducer, which are similar to those in Example

```haskell
mapper v = [(1, f v)] -- f stands for computation on input
reducer :: [(Integer, Integer)] -> (Integer, Integer)
reducer vs = (freq, total / freq)
  where (freq, total) = (sum (map fst vs), sum (map snd vs))
```

Then, \(\text{inv \ reducer}\) can be given by \(\{(\text{frq}, \text{avg}) \rightarrow \{(\text{frq}, \text{frq} \times \text{avg})\}\}\) and the generated functions work as expected.

It turns out that a sufficient condition for our synthesis to work is that a reducer is a homomorphism, defined below:

**Definition 1 (Homomorphisms [4])**. A function \(h\) of type \(\text{list} X \rightarrow Y\) is \(\odot\)-homomorphism for a binary operator \(\odot\) on type \(Y\) if and only if for arbitrary \(x\) and \(y\) of type \text{list} \(X\), the following holds:

\[
h(x + y) = h(x) \odot h(y).
\]

(Here, \(+\) stands for the function to append two lists.)

An intuitive reason why a list homomorphism works is that the value of applying the function does not depend on how a given list is split in calling the function recursively.

**Remark 2.** There are several options in the definition of a mapper, a combiner, and a reducer. For example, the MapReduce computation conducted by the following functions is also equivalent to that by \(m\) and \(r\).

```haskell
mapper’ :: data -> a
mapper’ v = m v

combiner’ :: [a] -> [a]
combiner’ vs = (inv r) (r vs)
```
The combiner defined this way outputs the return value of a right inverse. If \( r \) calculates the length and sum of a list, for example, the following \( \text{inv } r \) works as a right inverse for \( r \):

\[
(\text{inv } r)(s, n) = [s, 0, 0, \ldots, 0].
\]

However, this function \( \text{inv } r \) is less efficient than \text{mapper} \( , \text{combiner} \), and \text{reducer} because \( \text{inv } r \) yields the list whose length is equal to that of the input list; it does not improve the size of transferred data.

4 Correctness

This section formally states correctness of the method in Section 3. A combiner synthesized by our method is correct if it does not change the final result of the whole MapReduce computation in all the possible combiner-invocation scenarios. To formally state this property, we first define a computation model of MapReduce computation.

We call input to mappers \textit{input data}; we designate a meta-variable \( d \) for input data; we write \( D \) for the set of the input data.

4.1 Computation model

To motivate our computation model, we show a schematic presentation of a MapReduce computation in Figure 2. Input data, represented by a circle, are
fed to a mapper node; inside the mapper node, a mapper function converts the received data to a key–value pair represented by boxes. Then, a combiner aggregates the data. Those data are sent to the reducer node and processed by a reducer. In order to express this tree-like computation structure, we introduce computation trees.

**Definition 2 (Computation tree).** Computation trees, ranged over by a metavariable $t$, are defined by the following syntax.

$$ t ::= d \mid [t_1, t_2, \ldots, t_n]. $$

We write $T$ for the set of computation trees.

A computation tree is a rose tree where each leaf is labeled by an input datum. $[t_1, t_2, \ldots, t_n]$ represents one-step computation conducted by a combiner or a reducer called with the results of $t_1, \ldots, t_n$. A computation tree completely expresses when and to which data a combiner is applied. For example, the computation tree $["A 3", "A 5", "A 7"]$ represents the following computation: (1) the result of mapper applications to "A 3" and "A 5" are passed to a combiner and (2) the result of (1) and the result of the mapper application to "A 7" are passed to a reducer.

Suppose that a mapper $m'$, a combiner $c'$, and a reducer $r'$ are synthesized from a mapper $m$ and a list-homomorphic reducer $r$. Pick any input data $d_1, \ldots, d_n$ and any tree $t$ whose leaves collected from left to right are $d_1, \ldots, d_n$. Our method is correct if the result of the MapReduce computation expressed by $t$ with $m'$, $c'$, and $r'$ is equal to the one that is expressed by $[d_1, \ldots, d_n]$ with $m$ and $r$. This property can be formally stated by defining the result of the computation expressed by $t$. We define a function $\text{MR}_{m,c,r}$ to calculate the result in what follows.

We designate two auxiliary functions $\text{sort}$ and $\text{shuffle}$. They work as follows.

$\text{sort}$ sorts a list of key–value pairs by their keys. The ordering of keys are kept implicit here. To preserve the order of pairs with the same key, it has to be stable sorting. The type of $\text{sort}$ is $[(\text{Key}, \text{Value})] \to [(\text{Key}, \text{Value})]$. $\text{shuffle}$ takes a sorted list of key–value pairs as input and groups values by their keys; for example, $\text{shuffle} [("A", 3), ("A", 4), ("B", 5)]$ evaluates to $[("A", [3, 4]), ("B", [5])]$. The type of $\text{shuffle}$ is $[(\text{Key}, \text{Value})] \to [(\text{Key}, [\text{Value}])]$.

The concrete definitions of $\text{sort}$ and $\text{shuffle}$ are given in Appendix.

The following function $\text{ssf}$, taking a function $f$ and a list of key–value pairs $l$ as parameters, applies $\text{sort}$, $\text{shuffle}$, and $\text{map}$ $f$ to $l$ in this order.

**Definition 3 ($\text{ssf}$).**

$$ \text{ssf} ::= ([\text{Value}] \to \text{Value}) \to [(\text{Key}, \text{Value})] \to [(\text{Key}, \text{Value})] $$

$$ \text{ssf } f \ l = \text{map} (\text{id} \circ f) (\text{shuffle} (\text{sort} \ l)) $$

where $\text{id}$ is the identity function and $(p \triangle q)$ is the function defined by $(p \triangle q) x = (p x, q x)$. 

The function \( f \) passed to \( ssf \) is either a reducer or a combiner. \( ssf \) represents one-step aggregation conducted by a reducer or a combiner.

Using \( ssf \), the function \( MR_{m,c,r} \), which defines the result of the computation represented by a computation tree, is defined as follows.

**Definition 4 (\( MR_{m,c,r} \)).**

\[
\begin{align*}
MR_{m,c,r} : T &\rightarrow [(\text{Key}, \text{Value})] \\
MR_{m,c,r}^f x &= [m x] \\
MR_{m,c,r}^f [t_1, \ldots, t_n] &= ssf f (\text{concat} \ (\text{map} \ MR_{m,c,r}^c [t_1, \ldots, t_n])) \\
MR_{m,c,r} = MR_{m,c,r}^r.
\end{align*}
\]

The function \( MR_{m,c,r} \) \( t \) is the result of the MapReduce computation represented by \( t \) with a mapper \( m \), a combiner \( c \), and a reducer \( r \). The outermost level is aggregated by \( r \) whereas the inner level is aggregated by \( c \); this switching is represented by the auxiliary function \( MR_{m,c,r}^f \) where \( f \) represents the function used to aggregate the current level. If \( t \) is a tree of depth 1, we write \( MR_{m,c,r} t \) for \( MR_{m,c,r} \) \( t \); notice that a choice of a combiner does not affect the result of \( MR_{m,c,r} \) \( t \) in this case because \( c \) is never applied.

### 4.2 Soundness of the combiner

Fix a mapper \( m \) and a reducer \( r \). Let \( m' \), \( c' \), and \( r' \) be the mapper, the combiner, and the reducer generated by the program transformation in Section 3.

**Theorem 1 (Soundness of the combiner).** If \( r \) is a list homomorphism, then,

\[
MR_{m',c',r'} t = MR_{m-,r} (\text{collect} \ t)
\]

for any computation tree \( t \) where \( \text{collect} \) is defined as follows:

\[
\begin{align*}
\text{collect} x &= [x] \\
\text{collect} [t_1, \ldots, t_n] &= \text{concat} \ (\text{map} \ \text{collect} \ [t_1, \ldots, t_n]).
\end{align*}
\]

Note that \( \text{collect} \ t \) returns a tree of depth 1; hence, we use the notation \( MR_{m-,r} t \).

Theorem 1 follows from the following two lemmas.

**Lemma 1.** If \( r \) is a list homomorphism, then \( MR_{m',c',r'} t = MR_{m',c',r'} (\text{collect} \ t) \) for any computation tree \( t \).

**Lemma 2.** If \( r \) is a list homomorphism, then \( MR_{m',c',r'} (\text{collect} \ t) = MR_{m-,r} (\text{collect} \ t) \) for any computation tree \( t \).
5 Implementation

We implemented the algorithm in Section 3. Our system takes a mapper, a reducer, and its right inverse as input\(^5\) and generates Hadoop code equipped with a combiner. We assume that a reducer given to our system is specified with a language used by Morita et al. [9]. The language forces a reducer to be written in the following form:

\[
\begin{align*}
    f([a]) &= t_1; \\
    f([a] ++ x) &= t_2; \\
    f(x ++ [a]) &= t_3;
\end{align*}
\]

The first line is the definition of \(f\) for a singleton list. The second and the third line respectively defines the leftward and the rightward definitions of \(f\) for a list with two or more elements. By the third homomorphism theorem [4], there is a list homomorphism that conducts the computation specified by \(f\). In the current implementation, a programmer needs to ensure that the leftward version indeed defines the same function as the rightward version. The concrete syntax of the language is in the full version.

As mentioned above, the order of the original data at input is not preserved when they are passed to mapper nodes. Hence, a naive implementation of our combiner-synthesis method does not lead to a correct MapReduce program if the order of input data affects the result of computation. (Such order-sensitive computation is indeed frequently implemented with MapReduce code as reported by Xiao et al. [12].)

To cope with order-sensitive computation, our implementation generates code that manipulates key–value pairs extended with index ranges. An index range is a pair of integers \(\langle i_1, i_2 \rangle\) that represents physical positions from \(i_1\) to \(i_2\). We write \((k, v, \langle i_1, i_2 \rangle)\) for a key–value pair of \(k\) and \(v\) extended with an index range \(\langle i_1, i_2 \rangle\); Such a datum is obtained by assembling key–value pairs with the position from \(i_1\) to \(i_2\).

A mapper outputs key–value pairs; each pair is extended with a singleton index range \(\langle i, i \rangle\) where \(i\) represents the position of the generated pair at the input\(^6\). The combiner then assembles the data only if their index ranges are contiguous. The index part of the assembled data is merged; for example, a sum-computing combiner may assemble the list \([(A, 3, \langle 2, 2 \rangle), (A, 5, \langle 3, 3 \rangle), (A, 7, \langle 5, 5 \rangle)]\) to \([(A, 8, \langle 2, 3 \rangle), (A, 7, \langle 5, 5 \rangle)]\). Here, \((A, 3, \langle 2, 2 \rangle)\) and \((A, 5, \langle 3, 3 \rangle)\) may be assembled by the combiner because indices 2 and 3 are contiguous; the result then is \((A, 8, \langle 2, 3 \rangle))\) where \(2, 3\) represents a range of indices from 2 to 3.

One issue in the implementation is that triples passed to a reducer have to be sorted by indices in order-sensitive computation. However, Hadoop by default

---

\(^5\) The current implementation requires a user to manually write a right inverse of a reducer. Although Morita et al. [9] propose an automated right-inverse generation, their implementation is currently not available.

\(^6\) We assume that input data is preprocessed if necessary so that a mapper can calculate the appropriate index range for each pair.
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\[
\begin{array}{|c|c|c|c|}
\hline
\text{Benchmark} & \# \text{ of records} & \text{Data amount (MB)} & (\text{sec}) \\
\hline
\text{w/ combiner} & 1 \times 10^9 & 2.86 \times 10^{-3} & 120.5 \\
\text{w/o combiner} & 1 \times 10^8 & 6.98 \times 10^2 & 232.4 \\
\hline
\end{array}
\]

Table 1. The result of the Sum benchmark.

supports sorting only by keys. Our framework exploits secondary sorting \cite{11} to address this issue. Secondary sorting is a mechanism of Hadoop that enables sorting also by values, not only by keys. In the current implementation, the key–value pairs are sorted by indices and the values are grouped by the original single key. Hence, the resulting key–value pair contains values sorted by their indices.

These tricks about indices incur performance penalty also to order-insensitive problems, for which the order of input data does not matter. The current implementation provides an option to disable these tricks when a user knows her computation is an order-insensitive one. It is future work to automatically switch this option by deciding whether a problem is order-sensitive or not.

6 Experiment

We measured performance of synthesized code to confirm the effectiveness of our method. The experiment has been conducted on Amazon Elastic MapReduce \cite{1} with Hadoop 2.2.0. We used one master node and 10 worker nodes; they all have 7.5GB memory and 2 × 420 GB storage.

We measured performance of two benchmarks: Sum and MPS. We manually defined a right inverse of each function following the method proposed by Morita et al. \cite{9}. Sum calculates the sum of input data. MPS is the problem to compute the maximum sum of the prefixes of a given integer list. For example, the prefixes of \([4, -2, 7, -3]\) are \([4]\), \([4, -2]\), \([4, -2, 7]\) and \([4, -2, 7, -3]\); their sums are 4, 2, 9 and 6, respectively; hence, the MPS is 9.

The dataset used in the Sum benchmark consists of \(10^8\) 32-bit integers. We synthesized the code with the index-manipulating function mentioned in Section 5 disabled.

The datasets used in the MPS benchmark consist of \(10^9\) 32-bit integers associated with an integer index; the mapper parses the input data to get an integer \(n\) and an index \(i\), and generates \((\star, n, (i, i))\), where \(\star\) is a special designated symbol as the key. We used two datasets: the dataset Sequential, the integers are sorted by the indices; the integers in the dataset Random is ordered in random.

Table 1 shows the result of the Sum benchmark. We measured the number of communicated records between mapper nodes and reducer nodes, the amount of communicated data, and the time spent in the whole computation. The result suggests that our method reduces the amount of communicated data drastically in order-insensitive computation. This reduction contributes to speeding up the whole computation.
Problem Benchmark Sequential (# of records) Random (# of records)
MPS w/ combiner $1.0 \times 10^7$ $9.9 \times 10^7$
w/o combiner $1.0 \times 10^8$ $1.0 \times 10^8$

Table 2. The amount of transferred records.

<table>
<thead>
<tr>
<th>Problem</th>
<th>Benchmark</th>
<th>Sequential (MB)</th>
<th>Random (MB)</th>
</tr>
</thead>
<tbody>
<tr>
<td>MPS</td>
<td>w/ combiner</td>
<td>$4.64 \times 10^{-3}$</td>
<td>$2.06 \times 10^{3}$</td>
</tr>
</tbody>
</table>
w/o combiner | $1.40 \times 10^{3}$ | $1.41 \times 10^{3}$ |

Table 3. The amount of transferred data between mapper nodes and reducer nodes (1 MB = 1048576 bytes).

Table 2–4 show the result of the MPS benchmark. Table 2 presents the number of communicated key–value pairs between mapper nodes and reducer nodes, Table 3 the data amount, and Table 4 the time spent. For each dataset, we conducted experiment with the following programs: (1) a program with the generated combiner (the rows “w/ combiner”) and (2) a program without the generated combiner (the rows “w/o combiner”). The generated combiner works well in the dataset Sequential, whereas the combiner does not improve the performance in Random. This suggests that, if the data is preprocessed (i.e., sorted by indices) appropriately, our method works quite well also for order-sensitive problem. However, if this is not the case, the generated combiner incurs overhead.

7 Related work

Liu et al. [7] propose a transformation from sequential programs to MapReduce programs. In their framework, a program is specified by two functions: a (sequential) reducer and its right inverse written in Java. Then, their algorithm generates an equivalent MapReduce program via the third homomorphism theorem by Gibbons [4]. Actually, the way a new reducer is composed is the same as ours (although a right inverse is automatically generated by using a special language in our case), and so it could be used as a combiner. Nevertheless, their work does not consider combiners because the order of the input data is, in general, not preserved when the data are divided and sent to mapper nodes (personal communication with the authors). In our implementation, we use secondary sorting to solve the problem.

Morita et al. [9] develop a technique to derive weak right inverse and applied it to automatic parallelization based on the third homomorphism theorem, similarly to Liu et al. [7] above. We use their language and derivation of weak right inverse in our implementation. Our correctness theorem could be considered a variant of the third homomorphism theorem with additional twists caused by sorting and shuffling of key–value pairs. Although it is just complication, nei-
Problem | Benchmark | Sequential (sec) | Random (sec) |
--- | --- | --- | --- |
MPS | w/ combiner | 156.9 | 510.4 |
MPS | w/o combiner | 309.4 | 369.5 |

Table 4. Execution time. We measured the execution time with a Java method `System.currentTimeMillis()`.

ther Liu et al. [7] nor Morita et al. [9] deal with these additional complications rigorously.

Ono et al. [10] proposes a framework to verify correctness of a MapReduce program. They formalize a model of MapReduce computation with Coq [2] and proofs of its application-independent properties. Then, they provide a method for extracting Haskell code from Coq proofs. They also proposed another approach using Coq and Krakatoa [8], a verification tool for Java programs annotated with JML [6], in collaboration. The former approach can deal with only MapReduce programs such that a combiner is the same as a reducer. The latter approach could be applied to programs in which a combiner is different from a reducer. In both approaches, a programmer needs to provide a combiner implementation.

Yang et al. [13] also model MapReduce. They use CSP [5] as a tool for modeling, and mainly focus on the communication among mappers, reducers and the master node. In this model, however, a combiner is not considered.

8 Conclusion

We have proposed a framework for improving communication cost of MapReduce programs. Our method transforms a MapReduce program specified by a mapper and a reducer to an equivalent MapReduce program equipped with a combiner. The main idea is the use of a right inverse of the reducer. We have proved correctness of our synthesis. We have also implemented our method and confirmed its effectiveness with small experiment.

To apply our method, the input reducer has to be a homomorphism; this condition sometimes prohibits natural specification of a reducer, as we have seen in the example of an average-calculating program. Relaxing this restriction is one important direction of future work.

Our implementation, by using secondary sorting, deals with the fact that the order of input data may not be preserved when they are passed to mapper nodes. When the reducer is commutative (such as Sum), however, we do not have to apply secondary sorting. Optimization when the reducer is commutative is also interesting future work.

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Functional Kleene Closures

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Abstract. We present a derivation of a purely functional version of Kleene’s closure algorithm for Kleene algebras (with tests) that contain a subset where the closure is already known. In particular, our result is applicable to the Kleene algebra of square matrices over a given Kleene algebra. Our approach is based solely on laws imposed on Kleene algebras and Boolean algebras. We implement our results in the functional programming language Haskell for the case of square matrices and discuss a general implementation. In this process we incorporate purely algebraic improvements like the use of commutativity to obtain a concise and optimised functional program. Our overall focus is on a functional program and the computational structures from which it is composed. Finally, we discuss our result particularly in light of alternative approaches.

1 Introduction

The Kleene closure is a well established computational paradigm with numerous applications, e.g. in regular expressions. Also, it is possible to define the Kleene closure for matrices of Kleene algebras, which provides a unified approach to additional problems like the all-pair-shortest-path problem in graph theory or matrix multiplication or even inversion of matrices. Matrix algebras come with the additional benefit of being able to represent problems from other branches like reachability in graphs (Boolean matrices), the Dijkstra algorithm (matrices over the so-called tropical Kleene algebra, cf. Section 6) or the CYK algorithm (matrices over rules). This is to say that using matrices is not a restriction, but simply a general view for many different problems, while allowing all of the usual algebraic means associated with matrices.

The vast amount of problems captured by the computational scheme of this closure has led to a lot of research in this area, which includes several implementations in different programming languages. These are usually given as imperative (pseudo-)code, but there has been little development of a functional program. While there are several more or less canonical implementations in a functional programming language, they are usually based upon a translation of a given algorithm, but not on a purely functional approach. Clearly, there is a difference between a program in a functional language and a functional program, which is particularly important when dealing with algorithms. By definition an algorithm has a sequential look-and-feel that allows following a set of instruction
step by step. Such a construct fits well in the context of imperative languages, but may not be suited for a functional definition, since functional programs are usually inherently compositional and not necessarily computed in sequence.

In this paper we generalise an approach taken to compute a specific instance of the Kleene closure in [1] to the general case and present a purely functional program that can be used to compute the said closure. Our functions are prototypic by design, but the modularity of their components can be easily used to improve upon the implementation. This article is structured as follows.

- We recall the necessary preliminaries of a Kleene algebra and provide a definition of the Kleene closure that employs an auxiliary function.
- Using the algebraic reasoning we derive a recursive variant of this function.
- We implement the obtained recursion in Haskell, where we additionally employ Kleene algebra laws to improve performance.
- For comparison we implement the Kleene algorithm in two additional ways and perform tests.

To the best of our knowledge such a derivation has not been done so far. All of the presented code is given in Haskell [9]. In the course of the text we will refer to certain Haskell functions and modules all of which can be found using Hoogle (http://haskell.org/hoogle). A polished version of the code presented in this article is available at https://github.com/nikitaDanilenko/functionalKleene.

2 Algebraic Preliminaries

In the following we will deal with Kleene algebras according to the definition given in [7]. All definitions and consequences in this section are mentioned in the above source and the non-elementary results are proved as well; we include all of these for the sake of completeness only. We begin by the definition of a Kleene algebra, which can be split into two parts – the notion of an idempotent semiring and the concept of the star closure.

Definition 1 (Idempotent semiring and its order).
A structure \((S, +, \cdot, 0, 1)\) where \(+, \cdot\) are binary functions is called an idempotent semiring iff all of the following conditions hold:

(ISR1) \(+\) is associative and commutative with 0 as its neutral element
(ISR2) \(\cdot\) is associative with 1 as its neutral element and distributes over \(+\)
(ISR3) 0 is annihilating (i.e. \(\forall s \in S : 0 \cdot s = 0 = s \cdot 0\))
(ISR4) \(+\) is idempotent (i.e. \(\forall s \in S : s + s = s\))

The axioms (ISR1) and the first part of (ISR2) can be expressed more concisely: \((S, +, 0)\) is a commutative monoid and \((S, \cdot, 1)\) is a monoid. The explicit multiplication symbol \(\cdot\) is often omitted when there is no risk of ambiguity. Also, we always assume that multiplication binds tighter than addition. In an idempotent semiring we define for all \(s, t \in S\):

\[ s \leq t :\iff s + t = t. \]

Then \(\leq\) is an order and \(+, \cdot\) are monotonic (w.r.t. \(\leq\)) in both components.
With these preliminaries we are ready to define a Kleene algebra.

**Definition 2 (Kleene algebra).**

A structure \((K, +, \cdot, *, 0, 1)\), where \(+, \cdot\) are binary functions and \(*\) is a unary function, is called **Kleene algebra** (KA for short) iff all of the following hold:

(KA1) \((K, +, \cdot, *, 0, 1)\) is an idempotent semiring.

(KA2) \(\forall a \in K : 1 + a \cdot a^* \leq a^* \land 1 + a^* \cdot a \leq a^*\).

(KA3) \(\forall a, b \in K : (b \cdot a \leq b \rightarrow b \cdot a^* \leq b) \land (a \cdot b \leq a^* \cdot b \leq b)\).

Of the very numerous known properties of Kleene algebras we will use those that we summarise in the following lemma. All of these properties are proved in [7].

**Theorem 1 (Properties of \(*\) and the Kleene closure).**

Let \((K, +, \cdot, *, 0, 1)\) be a Kleene algebra. We define \(\star : K \rightarrow K\), \(a \mapsto a \cdot a^*\), called **Kleene closure**. Then the following hold:

(a) \(\forall a \in K : 1 + a \cdot a^* = a^*\).  
(b) \(\forall a, b \in K : 1 + a \cdot b = b \rightarrow a^* \leq b\).  
(c) \(\forall a, b \in K : (a + b)^* = a^* (b \cdot a^*)^*\).  
(d) \(\forall a, b, x \in K : a \cdot x = x \cdot b \rightarrow a^* \cdot x = x \cdot b^*\).  
(e) \(\forall a \in K : 1 + a^+ = a^* = (1 + a)^*\).

For the remainder of this section we assume that \((K, +, \cdot, *, 0, 1)\) is a Kleene algebra and \(n \in \mathbb{N}_{>0}\). Additionally, we will use the convention that \(*\) and \(\star\) bind tighter than multiplication. For many algebraic structures it is possible to lift the structure to the level of square matrices over the structure, which we denote by \(K_{n \times n}\). Kleene algebras are such a structure due to the following result. We write \(0_n\) and \(1_n\) to denote the zero and identity \(n \times n\)-matrices respectively. In the context of matrices we use "+" and "\(\cdot\)" to denote matrix addition and multiplication, respectively.

**Theorem 2 (\(K_{n \times n}\) as a Kleene algebra).**

\((K_{n \times n}, +, \cdot, *, 0_n, 1_n)\) is an idempotent semiring. Additionally, there is a function 
\(\star : K_{n \times n} \rightarrow K_{n \times n}\) such that \((K_{n \times n}, +, \cdot, *, 0_n, 1_n)\) is a Kleene algebra.

The first statement of this theorem is simple, but tedious to prove. The second is usually proved by defining an actual mapping \(\star : K_{n \times n} \rightarrow K_{n \times n}\). Throughout the literature one usually finds two particular definitions, which we mention here for the purpose of comparison. Let \(a \in K_{n \times n}\). For simplicity of notation we will index matrices starting from 0 instead of 1.

The first definition states that it is possible to compute matrices \(a^{(0)}, \ldots, a^{(n)}\) such that \(a^{(0)} := a\) and \(a^{(n)} = a^+\), where for all \(i, j, k \in \mathbb{N}_{<n}\) one has
\[
\begin{align*}
a^{(k+1)}_{i,j} := a^{(k)}_{i,j} + a^{(k)}_{i,k} \cdot \left(a^{(k)}_{k,j}\right)^* \cdot a^{(k)}_{k,j}.
\end{align*}
\]

This definition provides two ways of computing \(a^*\) – either as \(a^* = 1_n + a^+\) or as \(a^* = (1_n + a)^+\), by Theorem 1(e).
The second approach is based upon choosing \( l, m \in \mathbb{N}_{>0} \) such that \( n = l + m \) and splitting \( A \) into submatrices \( a = (p \, q) \) where \( p \in K^{l \times l} \), \( s \in K^{m \times m} \), and \( q, r \) have corresponding dimensions. Then one computes \( x := (p + q \cdot s^* \cdot r)^* \) and sets

\[
\begin{pmatrix}
oddspace x & x \cdot q \cdot s^* \\
\end{pmatrix}.
\]

(ii)

Since all of the matrices used in this definition have strictly smaller dimensions than \( a \) these computations can be used recursively for the computation of \( a^* \).

Both definitions are elegant in different ways – the first one is easy to translate in graph-theoretic terms and easy to implement in an imperative language, while the second one has a foundation in automata theory. Still, both definitions are rather algorithms (the second one can be even considered a non-deterministic one, since the actual decomposition does not matter), since they describe a sequence of steps that lead to \( a^* \). From a complexity point of view the second definition provides an additional challenge, since it contains matrix multiplication. It is not apparent at the first glance, what the exact complexity is, while the first one is clearly cubic in \( n \).

3 A Functional Approach

In this section we will develop a functional definition of the Kleene closure. Our approach is a direct generalisation of the methods used in [1]. For the sake of simplicity, but also for that of generality we need a slightly more sophisticated structure than a Kleene algebra, namely a so-called Kleene algebra with tests. Tests are elements of the Kleene algebra that “behave” like elements of a Boolean algebra. We use a similar notation and definition as in [8].

Definition 3 (Kleene algebra with tests).

A Kleene algebra with tests (KAT for short) is a structure \((K, B, +, \cdot, *, 0, 1)\) such that all of the following conditions hold:

(KAT1) \((K, +, \cdot, *, 0, 1)\) is a Kleene algebra.

(KAT2) \(B \subseteq K\).

(KAT3) \((B, +_B, \cdot_B, 0, 1)\) is a Boolean algebra (i.e. a distributive and complementory lattice), where \(+_B, \cdot_B\) are the operations of \(K\) restricted to \(B\).

We suppress the explicit negation operation for simplicity, since it is easily recovered from the complementarity that provides complements and distributivity, which can be used to show that complements are unique. Note that for any Kleene algebra there is a trivial set of tests namely \(B_0 = \{0, 1\}\). Abstract tests provide an elegant means to express Boolean conditions without leaving the Kleene algebra. All tests \(b \in B\) satisfy \(b \leq 1\), because in Boolean algebras 1 is the largest element and thus \(b + 1 = 1\), which yields \(b \leq 1\) in \(K\).

From now on we assume that \((K, B, +, \cdot, *, 0, 1)\) is a Kleene algebra with tests. We then consider the following function.
\[ \tau : K \times B \rightarrow K, \quad (a, b) \mapsto a \cdot (b \cdot a)^* \].

This function is a translation of the function in relational terms from [1]. Just as its relational version, \( \tau \) has the following properties for all \( a \in K \):

\[
\begin{align*}
\tau(a, 0) &= a \cdot (0 \cdot a)^* = a \cdot (0^*) = a \cdot 1 = a, \\
\tau(a, 1) &= a \cdot (1 \cdot a)^* = aa^* = a^+.
\end{align*}
\]

(iii)

To deal with tests between 0 and 1 we take a similar approach as in [1] and study the recursion properties of \( \tau \). We observe that for all \( a \in K \) and all \( b_1, b_2 \in B \) we get the following chain of equations:

\[
\begin{align*}
\tau(a, b_1 + b_2) &= a \cdot ((b_1 + b_2) \cdot a)^* \\
&= a \cdot (b_1 a + b_2 a)^* \\
&= a \left( (b_1 a)^* (b_2 a) (b_1 a)^* \right)^* \\
&= (a (b_1 a)^*) (b_2 (a (b_1 a)^*))^* \\
&= \tau(a, b_1) \cdot (b_2 \cdot \tau(a, b_1))^* \\
&= \tau(\tau(a, b_1), b_2) \\
&= \tau(a, b_1 + b_2) = \tau(a, b_1) (b_2 \tau(a, b_1))^* = \tau(\tau(a, b_1), b_2). \quad \text{(iv)}
\end{align*}
\]

This property is a generalisation of the one derived in [1], both in terms of the decomposition of the underlying set as well as the algebraic structure at hand, but the steps on the way are exactly those used in the above source.

Now consider the application of the above recursion to a finite subset \( B' \subseteq B \), such that\(^1\) \( B' = \{ b_0, \ldots, b_m \} \) for some \( m \in \mathbb{N} \) and

\[ \sum_{i=1}^{m} b_i = 1. \]

Then we can use the above formula and compute \( a_0 := a \) and \( a_{i+1} := \tau(a_i, b_{i+1}) \) for all \( i \in \{ 1, \ldots, m \} \) to obtain \( a_m = a^+ \). Note that this is very similar to the construction given in the introduction only that it does not depend on matrices.

This computational paradigm is captured by a left-fold. Note that it does not depend on a particular order of traversal of the partition elements, but is intrinsic to the actual computation. In the above example we split the sum into its first summand and the rest, while we could have taken the last element and the corresponding rest as well without changing the result. Left-folds are favoured over right-folds in strict languages since they are usually more efficient (tail-call). In a lazy setting tail-calls can become more complex, because the accumulation

\(^1\) Such a set where \( b_i < 1 \) for all \( i \in \mathbb{N}_{<m} \) does not necessarily exist. However, if \( B \) is finite, we can simply choose the atoms of \( B \) (i.e. the upper neighbours of 0) as \( B' \).
parameter is not evaluated until needed, while its construction grows increasingly more complex. Since we are looking for a solution in a lazy functional language, we may need to transform the above recursion into a (non-generic) right-fold. Our aim is to use the function $\tau$ for the computation of the Kleene closure of a given algebra $K$. To do that we will determine the Kleene closure for a specific subset of $K$, namely $\{ ba \mid b \in B, a \in K \}$ and then apply $\tau$ in a recursive fashion as described above, assuming that there is a decomposition of 1 as a sum of finitely many tests. By Equation (iv) the knowledge of $^*$ on the above set is enough to compute the Kleene closure of any element of the Kleene algebra.

Let us summarise the result of this section in the following theorem.

**Theorem 3 (Recursive computation of $\tau$).**

Let $(K, B, +, \cdot, ^*, 0, 1)$ be a KAT and $\tau : K \times B \to K$, $(a, b) \mapsto a(ba)^*$. Then the following hold:

1. For all $a \in K$ and $b, c \in B$ we get $\tau(a, b + c) = \tau(a, b)(\tau(a, b))^*$.
2. For all $m \in \mathbb{N}$ and $b \in B^m$ such that $\sum_{i=0}^{m-1} b_i = 1$, all $a \in K$, all $i \in \mathbb{N}_{<m}$ setting $n_i := \sum_{j=0, j \neq i} b_j$ we find that $a^+ = \tau(a, 1) = \tau(a, n_i) (b_i \tau(a, n_i))^*$.

**Proof.** The first claim is just a rephrasal of Equation (iv). For the second let $m \in \mathbb{N}$, $b \in B^m$, $a \in K$, $i \in \mathbb{N}_{<m}$ and $n_i$ as required. Then $1 = n_i + b_i$ and thus

$$a^+ \overset{\text{Eq. (iii)}}{=} \tau(a, 1) = \tau(a, n_i + b_i) \overset{(1)}{=} \tau(a, n_i) (b_i \tau(a, n_i))^* .$$

In statement (ii) of the previous theorem the $b_i$ is removed from the sum for clarity, but the theorem obviously holds without this removal as well, because in idempotent semirings (e.g. the Boolean algebra $B$) addition is idempotent. As for applications of the last equality, the motivation is that a good choice of $b$ allows a simple computation of $(b_i \tau(a, n_i))^*$, which can then in turn be applied to compute $a^+$ or $a^*$ by iteration.

### 4 Application to Square Matrices

In this section we apply the technique from the previous section to the Kleene algebra of square matrices with entries from a Kleene algebra. We use the Boolean algebra of partial identities as test, which we elaborate shortly. For the remainder of this section let $K$ be a Kleene algebra and $n \in \mathbb{N}$. Additionally, we do not differentiate between $K^{1 \times 1}$ and $K$, just as we consider the sets $K^n$ and $K^{1 \times n}$ to be the same. To be perfectly accurate we should use isomorphisms between these sets, but we omit these to avoid unnecessary clutter.

To avoid confusion, we will index constants by their algebra throughout this section and make ample use of brackets to avoid indexing the star operation. Recall that $1_n$ is simply the identity matrix and $0_n$ is the zero matrix. For every $a \in K^{n \times n}$ we use $a_j$ to address the $j$-th row of $a$ and $a_{j,k}$ to address the component at the $j$-th row in the $k$-th column.

We define partial identities by using standard unit vectors.
Definition 4 (Standard unit vectors and matrices, partial identity). We define for all $i \in \mathbb{N}_{<n}$ the $i$-th standard unit row vector $e_i \in K^n$ as 
$$e_i : \mathbb{N}_{<n} \to K, \ p \mapsto \delta_{i,p},$$
where $\delta$ is the Kronecker-delta. Clearly, $e_i \cdot e_j^\top = \delta_{i,j}$ for all $i, j \in \mathbb{N}_{<n}$. Additionally for all $i, j \in \mathbb{N}_{<n}$ we define the $(i, j)$-th standard unit matrix as $e(i, j) := e_i^\top \cdot e_j$. For every $S \subseteq \mathbb{N}_{<n}$ the partial identity on $S$ is defined as 
$$I_S := \sum_{i \in S} e(i, i)$$

and $\Pi(K, n) := \{ I_T \mid T \subseteq \mathbb{N}_{<n} \}$.

As mentioned above, we use the set of all partial identities $\Pi(K, n)$ as tests, which requires this set with the restricted operations of $K^{n \times n}$ to be a Boolean algebra. This is indeed true.

Lemma 1 (Partial identities are a Boolean algebra). Let $B := \Pi(K, n)$. Then the structure $(B, +_B, \cdot_B, 0_n, 1_n)$ is a Boolean algebra, where $+_B, \cdot_B$ are the operations of $K^{n \times n}$ restricted to $B$.

Proof. We provide only an outline of the proof. Consider the function 
$$f : (2^{\mathbb{N}_{<n}}, \cup, \cap, \emptyset, \mathbb{N}_{<n}) \to (B, +_B, \cdot_B, 0_n, 1_n), \ S \mapsto I_S.$$  
Clearly, $f$ is bijective, $f(\emptyset) = 0_n$ and $f(\mathbb{N}_{<n}) = 1_n$. A simple, but slightly lengthy, computation yields that $f$ distributes over addition and multiplication as well. Thus $f$ is a constant preserving lattice isomorphism. Since its domain is a Boolean algebra, so is its image, which is its range, because $f$ is bijective. \qed

To apply the function $f$ from the previous section, we need a finite decomposition of $1_n$ into tests. We choose a decomposition in standard unit matrices by setting $b_i := e(i, i)$ for all $i \in \mathbb{N}_{<n}$. As stated just after Theorem 3 this choice should provide a simple means to compute $(b_i a)^\ast$ for all $i \in \mathbb{N}_{<n}$ and all $a \in K^{n \times n}$. To see that this is the case, we need two auxiliary lemmas.

Lemma 2 (Homothetic injection). Let $\varphi : K \to K^{n \times n}$, $c \mapsto c \cdot 1_n$, where $\cdot : K \times K^{n \times n} \to K^{n \times n}$ is the scalar multiplication of a matrix. Then $\varphi$ is a Kleene algebra homomorphism.

Proof. Clearly, $\varphi(0K) = 0_n$ and $\varphi(1K) = 1_n$ by the very definition of scalar multiplication. The additivity and multiplicativity of $\varphi$ are immediate consequences of this definition as well. The only difficulty is the star operation. Let $a \in K$. Then we have $1_n + \varphi(a) \cdot \varphi(a)^\ast = \varphi(1K) + \varphi(a \cdot a^\ast) = \varphi(1K + a a^\ast) = \varphi(a^\ast)$ by Theorem 1.(a). By Theorem 1.(b) this provides $\varphi(a)^\ast \leq \varphi(a^\ast)$. Thus $\varphi(a)^\ast$ has non-zero entries at most along its diagonal. Let $i \in \mathbb{N}_{<n}$. Then we have:
$$\begin{align*}
(\varphi(a)^\ast)_{i,i} &= (1_n + \varphi(a)\varphi(a)^\ast)_{i,i} = 1_K + ((a \cdot 1_n) \varphi(a)^\ast)_{i,i} \\
&= 1_K + (a \cdot \varphi(a)^\ast)_{i,i} = 1_K + a (\varphi(a)^\ast)_{i,i},
\end{align*}$$
which again by Theorem 1.(b) yields $\varphi(a)^\ast_{i,i} = a^\ast \leq (\varphi(a)^\ast)_{i,i}$. This results in $\varphi(a)^\ast \leq \varphi(a)^\ast$ and since $\leq$ is an order, we finally conclude $\varphi(a)^\ast = \varphi(a)^\ast$. \qed
Lemma 3. Let \( c, v \in K \) such that \( v^2 = cv \). Then \( v^* = 1 + c^* v \).

Proof. Set \( a := c, b := v \) and \( x := v \). Then \( ax = cv = vv = xb \), hence by Theorem 1.(d) we get \( c^* v = a^* x = xb^* = vv^* \), which yields \( v^* = 1 + vv^* = 1 + c^* v \). □

With this lemma we can compute \((e(i,i)a)^*\) without using the star in \( K^{n \times n} \).

Theorem 4 (Kleene closure of matrices).

Let \( a \in K^{n \times n} \), \( b \in \Pi(K, n) \) and \( i \in \mathbb{N}_{\leq n} \). Then we have:

1. For all \( x \in K^{n \times n} \) and \( k \in \mathbb{N}_{\leq n} \) we have \((e(k,k)x)^* = 1_n + ((x_{i,i})^*) \bullet e(k,k)x \).

2. \( \tau(a, b + e(i, i)) = \tau(a, b) + \tau(a, b) (((\tau(a, b)_{i,i})^*) \bullet e(i,i)) \tau(a, b) \).

3. For every \( j \in \mathbb{N}_{\leq n} \) we additionally get

\[
\tau(a, b + e(i, i))_j = \tau(a, b)_j + \tau'(\tau(a, b)_{j,i} \cdot (\tau(a, b)_{i,i})^*) \tau(a, b)_i ,
\]

where \( \tau' \) is the addition and \( \tau \) is the scalar multiplication of vectors.

Proof. (1) Let \( x \in K^{n \times n} \), \( k \in \mathbb{N}_{\leq n} \) and \( \varphi \) as in Lemma 2. Then we get:

\[
e(k,k)x \cdot e(k,k)x = e_k^T e_k x e_k^T e_k x = e_k^T (x_k,k) e_k x = x_k,k \bullet (e_k^T e_k x) = \varphi(x_k,k) \cdot e(k,k)x .
\]

By Lemmas 3 and 2 this yields that

\[
(e(k,k)x)^* = 1_n + \varphi(x_k,k)^* e(k,k)x = 1_n + \varphi((x_k,k)^*) e(k,k)x = 1_n + (x_k,k)^* \bullet e(k,k)x .
\]

(2) We calculate as follows:

\[
\tau(a, b + e(i, i)) = \tau(a, b)(e(i,i)\tau(a,b))^* \quad \text{by Theorem 3}
\]

\[
= \tau(a, b) (1_n + (\tau(a,b)_{i,i})^* \bullet e(i,i)\tau(a,b)) \quad \text{by (1)}
\]

\[
= \tau(a, b) + \tau(a, b) (\tau(a,b)_{i,i})^* \bullet e(i,i) \tau(a, b) .
\]

(3) First of all for every \( d \in K^{n \times n} \), \( c \in K \) and \( j \in \mathbb{N}_{\leq n} \) the following holds

\[
(d \circ e(i,i) d)_j = e_j d (c \circ e_i^T e_i d) = e_j d e_i^T (c \circ e_i d) = (d_{j,i} c) * d_i . \quad (!)
\]

Let \( j \in \mathbb{N}_{\leq n} \). Then the following holds:

\[
\tau(a, b + e(i, i))_j = (\tau(a, b) + \tau(a, b) (\tau(a,b)_{i,i})^* \bullet e(i,i) \tau(a,b))_j \quad \text{by (2)}
\]

\[
= \tau(a, b)_j + \tau'(\tau(a, b)_{j,i} \cdot (\tau(a, b)_{i,i})^*) \tau(a, b)_i \quad \text{by (1) □}
\]
5 A Functional Implementation

In this section we use the approach of the previous sections to obtain a functional implementation for Kleene closure of a square matrix. We discuss this restriction in a moment. The first step in the implementation is to encode the required algebraic structures. One particularly simple way to do that is to use type classes.

```haskell
class IdempotentSemiring σ where
  (⊕), (⊙) :: σ → σ → σ
  zero, one :: σ
  isZero, isInOne :: σ → Bool

class IdempotentSemiring κ ⇒ KleeneAlgebra κ where
  star :: κ → κ
```

Additionally, we require instances of these type classes to satisfy the corresponding algebraic laws. Such an approach requires a user to check the necessary conditions, which may or may not be neglected in an application\(^2\). The predicates `isZero`, `isOne` are not explicit parts of the algebraic definition. In a purely theoretical context we can always compare values for equality, while in practice the equality of two objects may be undecidable. In Haskell the `Eq` type class represents values that support equality checks and we could have used this type class as a superclass for `IdempotentSemiring`. The reason we did not is that the above approach is more general and checking the constants is sufficient to obtain some simple improvement. Further conditions that may arise for specific Kleene algebras may be added to more specialised functions dealing with these algebras.

For a general setting, we can define a type class `KAT` for Kleene algebras with tests as a subclass of `IdempotentSemiring` that provides a constant `sparseTests :: KAT κ ⇒ [κ]` which denotes a set of tests that is minimal with respect to the fact that the sum of these tests is `one`. Then we can use Theorem 3 as follows:

```haskell
katStar :: KAT κ ⇒ κ → κ
katStar a = one ⊕ katPlus a

katPlus :: KAT κ ⇒ κ → κ
katPlus a = tau a sparseTests

tau :: KAT κ ⇒ κ → [κ] → κ

tau a [] = a

tau a (t : ts) = x ⊕ katStar (t ⊕ x) where x = tau a ts
```

However, to use this implementation for square matrices one has to provide a semiring instance for said matrices. This in turn requires taking the matrix sizes into account, which comes with additional overhead (e.g. checking sizes by hand or encoding the sizes in the matrix type). For the sake of demonstration we take a different approach in the case of square matrices using the specialised Theorem 4 for a prototypical implementation without a semiring instance for matrices.

\(^2\) To ensure that the required properties are in fact satisfied one can use automated tools like Coq (cf. [2]) to make certain functions applicable only once the precondi-
Theorem 4 provides us with a possible implementation assuming certain existing functions. Both equations given in said theorem require access to single values of a matrix, i.e. $A_{i,j}$. Thus if we want to use these equations for an implementation we need some representation of matrices that allows such queries. Also, the first equation uses matrix multiplication while the second one does not, but instead relies on the concept of rows. There are numerous ways for a representation of a matrix, particularly if it is sparse (i.e. small percentage of non-zero values) and different representations provide different features and caveats. For now we choose the adjacency list model as a middle ground between efficiency and simplicity. To that end we use the following notations.

$$\text{type } \text{Row } \alpha = [(\text{Int}, \alpha)]$$
$$\text{type } \text{Mat } \alpha = [\text{Row } \alpha]$$

$$(\cdot): \text{IdempotentSemiring } \sigma \Rightarrow \text{Row } \sigma \rightarrow \text{Int } \rightarrow \sigma$$
$$\sum = \text{zero}$$
$$((i, v): \text{ivs}) \cdot k \mid i \equiv k = v$$
$$i < k = \text{ivs} \cdot k$$
$$\text{otherwise} = \text{zero}$$

Additionally we assume and maintain the conditions that Rows are sorted in ascending order of their first components, that the second components are non-zero and that an $n \times n$-matrix $A$ is represented by a list of rows $a$ such that $A_i = a!!i$ and $A_{i,j} = (a!!i) \cdot j$ for all $i, j \in \mathbb{N}_{<\alpha}$, where (!!) is the built-in Haskell function for accessing indices of a list. This representation is similar to the one in [1, 3] and is provided for the completeness of the implementation.

To apply Theorem 4 we need a scalar multiplication and addition of rows. A straightforward implementation of the former is obtained by multiplying every second component of a row with a given scalar and tidying up the result to remove possible zeroes. There is room for some canonic improvement, since scalar multiplication in case of 0,1 is particularly simple. To that end we can simply define the actual multiplication as follows.

$$(\cdot): \text{IdempotentSemiring } \sigma \Rightarrow \sigma \rightarrow \text{Row } \sigma \rightarrow \text{Row } \sigma$$

$$s \cdot \text{row} \mid \text{isZero } s = []$$
$$s \cdot \text{row} \mid \text{isOne } s = \text{row}$$
$$\text{otherwise} = \text{filter } (\neg \circ \text{isZero } \circ \text{snd}) \cdot (\text{map } (\lambda(i, v) \rightarrow (i, s \circ v)) \cdot \text{row})$$

For all $x, y \in K$ with $x+y = 0$ we get $x = x+0 = x+(x+y) = (x+x)+y = x+y = 0$, which results in $y = 0 + y = x + y = 0$ and thus non-zero values are closed under addition. That is to say that in an idempotent semiring we don’t need the tidying step $\text{filter } (\neg \circ \text{isZero } \circ \text{snd})$ when adding vectors. Our precondition that the first components of the rows are increasingly sorted allows implementing the addition of rows in terms of a straightforward merging strategy.

---

3 In Kleene algebras there can be zero divisors, i.e. elements $x, y \in K$ such that $x \neq 0$ and $y \neq 0$, but $xy = 0$. This is why we need to filter the zero values.
\[ (x') :: \text{IdempotentSemiring} \rightarrow \text{Row} \rightarrow \text{Row} \rightarrow \text{Row} \sigma \]

\[
\begin{array}{c|c|c}
\tau & x & {[ y +' x ]} \\
\hline
{} & y & y \\
\hline
x & x & x \\
\end{array}
\]

\[ x(\{(i, v) : \text{ivs}\} +' y(\{(j, w) : \text{jws}\} | i \equiv j = (i, v \oplus w) : \text{ivs} +' \text{jws}) | i < j = (i, v) : (\text{ivs} +' y) | \text{otherwise} = (j, w) : (x +' \text{jws}) \]

The function \( \tau \) which we want to use for the computation of the Kleene closure depends on a finite decomposition of 1 into tests. In our approach in the previous section we used the tests \( \{ e(i, i) | i \in \mathbb{N}_{<n} \} \). Every such test and every sum of such tests is uniquely determined by a subset of \( \mathbb{N}_{<n} \), which we have established in the proof of Lemma 1. Thus instead of successively computing \( b_i \) and \( n_i \) as in Theorem 3 we can use the decomposition of tests in \( (\mathbb{N}_{<n} \setminus \{ i \}) \cup \{ i \} \) instead of the corresponding test. To gather the indices of a matrix we zip the matrix with the list of natural numbers and ignoring the matrix values.

\[ \text{spine} :: [\alpha] \rightarrow [\text{Int}] \]

\[ \text{spine} = \text{zipWith const} [0..] \]

The Kleene closure of \( a \) is simply \( \tau(a, f(\mathbb{N}_{<n})) \) for every \( a \in K^{n \times n} \), where \( f \) is the isomorphism between tests and sets from Lemma 1. Disregarding the \( f \) and using lists to represent sets we can implement this in Haskell as follows.

\[ \text{kleeneClosure} :: \text{KleeneAlgebra} \kappa \Rightarrow \text{Mat} \kappa \rightarrow \text{Mat} \kappa \]

\[ \text{kleeneClosure} a = \tau(a) (\text{spine} a) \]

The function \( \tau \) can then be defined in an inductive way. First the base case.

\[ \tau :: \text{KleeneAlgebra} \kappa \Rightarrow \text{Mat} \kappa \rightarrow [\text{Int}] \rightarrow \text{Mat} \kappa \]

\[ \tau a [\] = a \]

This is simply a quotation of a property of \( \tau \). Finally, \( \tau(a, f(\{ k \} \cup S)) \) is a matrix whose rows are given by Theorem 4.(3). We can then write

\[ \tau a (k : s) = \text{newMat} k (\tau a s) \]

and thus delay the actual computation in the auxiliary function \( \text{newMat} \) that captures the computation scheme of the above proposition.

\[ \text{newMat} :: \text{KleeneAlgebra} \kappa \Rightarrow \text{Int} \rightarrow \text{Mat} \kappa \rightarrow \text{Mat} \kappa \]

\[ \text{newMat} i a = \text{map} (\lambda a_j \rightarrow (a_j ! i) \odot \text{star} (a_i ! i) \ast a_i +' a_j) a \text{ where } a_i = a!!i \]

In essence this definition is the implementation of the equation from Theorem 4.(3). The only difference is that we swapped the arguments of \( (+') \), which is a valid transformation, since \( + \) is commutative, which makes \( (+') \) commutative as well. Experiments have shown that this simple algebraic rule improves the running times by a factor between 1.5 and 2. This concludes the implementation.

We observe that the full definition of \( \tau \) is

\[ \tau a [\] = a \]

\[ \tau a (k : s) = \text{newMat} k (\tau a s) \]

which by the universal property of \( \text{foldr} \) (cf. [5]) yields \( \tau a \equiv \text{foldr} \text{newMat} a \) that in turn can be \( \eta \)-reduced to \( \tau \equiv \text{foldr} \text{newMat} \).
6 Alternative Implementations and Comparison

In this section we discuss two slightly more canonical approaches to an implementation of the Kleene closure function and test their complexities in terms of running times and space consumptions.

6.1 Using Arrays

As we have stated before, Equation (i) provides an out-of-the-box algorithm in imperative languages using arrays. Haskell provides a number of different arrays (mutable and immutable) with fast access to their indices. For the sake of simplicity we used the most basic arrays that allow simple and pure code.

The downside of these arrays is that modification is very costly. Fortunately, in our application we don’t need to modify arrays but only build new ones. Also, by Equation (i) we don’t require \( n \) different arrays, but only two different ones, since \( a^{(k+1)} \) is constructed from \( a^{(k)} \) alone. This is to say that while we compute \( n \) different arrays, only two of them need to be present in memory at any given time, namely the \( k \)-th and \( 1+k \)-th ones, while all previous versions are not required and thus garbage collected. When we represent matrices by arrays in the canonical fashion, we may implement the Kleene closure as follows.

```
newtype ArrayMat α = ArrMat { unArrMat :: Array (Int, Int) α }
kleeneClosureArray :: KleeneAlgebra κ ⇒ ArrayMat κ → ArrayMat κ
kleeneClosureArray (ArrMat a₀) = foldl newArray a₀ [0..n] where
  newArray arr k = listArray bnds (map (newValue arr k) positions)
  bnds = bounds a
  positions = range bnds
  n = snd (snd bnds)
  newValue a k (i,j) = (a! (i,k)) ⊙ star (a! (k,k)) * a! (i,j)
```

Here (!) is the array query function and not our row querying function from the previous section. The local function newValue produces the value \( a^{(k+1)}_{i,j} \), newArray maps this producer over all index pairs \((i,j)\) and foldl repeats this procedure for all \( k \) and basically acts as a for-loop.

6.2 Another List Version

Another version comes to mind by simply observing that Equation (i) contains the index \( j \) at the same relative position. This allows the following rephrasal:

\[
\forall i,k \in \mathbb{N} : a^{(k+1)}_i = a^k_i + a^k_i \cdot (a^k_{i,k})^* \cdot a^k_k.
\]

This is strikingly similar to our recursion for \( \tau \). The essential difference is that the values of \( k \) are traversed in another direction and it is not quite obvious why this is a valid transformation. The specification of Equation (v) can be used to obtain a third implementation, where newMat is the function from Section 5.

```
kleeneClosureLeft :: KleeneAlgebra κ ⇒ Mat κ → Mat κ
kleeneClosureLeft a = foldl (flip newMat) a (shape a)
```
6.3 Comparison

We have implemented all three closure functions in essentially the presented way. Additionally, we implemented a random matrix generation based upon shuffling. Given a density \( d \in [0, 1] \) and a size \( n \) we compute the number of non-zero positions as \( p = \lfloor d \cdot n^2 \rfloor \). Then the first \( p \) positions in the matrix are filled and then the matrix is shuffled. This technique is known to be uniformly distributed, is already implemented in Haskell and it depends only on a single random generator, which can be created using a single \texttt{Int}. Additionally, random generators in Haskell generate pseudo-random numbers due to referential transparency, i.e., taking a random number from a fixed random generator will always produce the same result. This makes testing of randomly generated data easily repeatable. We randomly generated three random number generators from the random number generator generated from the number 42 and ran our three implementations on different sizes \( n \), densities\(^{5}\) \( d \) and Kleene algebras. In all cases the generation of the random matrix counts towards the total time to simulate pre-processed input. The result of every function is fully evaluated, but not printed. The values in Tables 1 and 2 display the average (avg) and the maximum (max) space consumption in megabyte and the running time (sec) in seconds and all values are arithmetic means over all generators\(^{6}\).

We observe that the developed right-fold version is almost always better than both other versions in terms of space and time consumption and if it’s not, it is very close to the left-fold variant. The reason for the early failure in the Boolean semiring is Haskell’s non-strictness, which in this case needs to be tamed with non-algebraic means. Still, the right-fold function manages all matrices of size 750, while both other versions do not. One typical improvement of tail-recursive functions in Haskell is the use of so-called strictness annotations which

\(^{4}\) The tropical KA is \((\mathbb{N} \cup \{ \infty \}, \min, +, \infty, 0, \star)\), where \(\star\) is the constant 1-function.

\(^{5}\) In a matrix with size 1000 a density of 0.1 means 100000 entries, which is likely to yield a fully filled transitive closure. This is why we use such seemingly small values.

\(^{6}\) The measurements were taken on a machine with an Intel Core i5-2520M CPU (4 \( \times \) 2.5 GHz) with 8 GB of DDR3 RAM, running Ubuntu 12.04 and using GHC 7.6.3.
(partially) evaluate a parameter before usage. We have experimented with this technique in the above implementations and found that it yields little to none improvement in case of the left-fold variant, but considerable (yet constant) improvements in the right-fold version. We take this as an indicator that the computational paradigm of a right-fold is conceptually better suited for the Kleene closure than that of a left-fold.

7 Related Work and Discussion

There exist different approaches to a functional version of Kleene's algorithm. In [10] an implementation is given that is based upon Equation (i) and the use of arrays and the more recent work [4] uses Equation (ii) for an implementation. The special case of transitive closures (closures over the Boolean KA) is treated in [6] using a monadic abstraction of lazy arrays that is implemented efficiently internally. While these articles are very well-written, they do not feature an actual derivation of the implementation and depend on the chosen data types. In our implementation we made an implementation choice, but only to supply the complete code. It is very simple to abstract our implementation to a general one – in fact, the only place we explicitly use the fact that matrices are represented by lists of rows is the implementation of \texttt{newMat} that requires a (!!!). Obviously, this implementation can be replaced by a more parametric one, which is parametrised over the container type for matrices, and then also over the one for rows. In conclusion, any representation of matrices that supports the notion of rows and rows that support addition, query and scalar multiplication is suited for our implementation. We used KATs to derive a general version of a Kleene closure. While KATs are less general than simple Kleene algebras, the Kleene algebra of square matrices over a Kleene algebra is a KAT without additional preconditions, because the set of partial identities is a Boolean algebra (cf. Lemma 1). Thus our approach for square matrices is applicable in the same settings as mentioned above and the one we took before the specialisation is applicable to all KATs where 1 is decomposable into a (finite) set of tests. Note that the right-fold structure can yield partial values in the infinite case as well, if some information
about \( x + y \) can be extracted from \( x \) alone. This is a proper improvement of the other versions from the Section 6, because left-folds diverge on infinite lists.

We have dealt with a single algorithm in this article only, but the employed techniques can transferred to more problems, too. For instance, a very similar approach to a graph representation has been taken in [3] to compute maximum matchings. On a more general note it is natural to take a row-based approach, since the rows of a matrix represent the successor lists of a given graph and it is likely that algorithms that can be elegantly expressed via successor lists have an elegant algebraic representation, too. In the beginning we mentioned that algorithms that are expressed in terms of a Kleene closure in some specific algebra can often be rewritten as a matrix closure. In fact every such closure can be expressed trivially as a matrix closure by Lemma 2.

We are confident that many more functional implementations can be obtained from algebraic specifications, especially in the field of graph algorithms. Additionally, the combination of algebraic reasoning and functional programming can reveal complexity bottlenecks through parametric abstraction over the structures. We have shown that it is possible to find a specification that is based upon modifying a value that will be computed next instead of one using values that have already been computed. While the latter is tail-recursive and can yield a performance gain in a strict setting, the former is better suited for a non-strict approach allowing propagation of partial values before the recursive application.

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**References**

On completeness of logic programs

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Abstract. Program correctness (in imperative and functional programming) splits in logic programming into correctness and completeness. Completeness means that a program produces all the answers required by its specification. Little work has been devoted to reasoning about completeness. This paper presents a few sufficient conditions for completeness of definite programs. We also study preserving completeness under some cases of pruning of SLD-trees (e.g. due to using the cut). We treat logic programming as a declarative paradigm, abstracting from any operational semantics as far as possible. We argue that the proposed methods are simple enough to be applied, possibly at an informal level, in practical Prolog programming. We point out importance of approximate specifications.

Keywords: logic programming, program completeness, declarative programming, approximate specification.

1 Introduction

The notion of partial program correctness splits in logic programming into correctness and completeness. Correctness means that all answers of the program are compatible with the specification, completeness – that the program produces all the answers required by the specification.

In this paper we consider definite clause programs, and present a few sufficient conditions for their completeness. We also discuss preserving completeness under pruning of SLD-trees (by e.g. using the cut). We are interested in declarative reasoning, i.e. abstracting from any operational semantics, and treating program clauses as logical formulae. Our goal is simple methods, which may be applied – possibly informally – in actual practical programming.

Related work. Surprisingly little work was devoted to proving completeness of programs. Hogger [10] defines the notion of completeness, but does not provide any sufficient conditions. Completeness is not discussed in the important monograph [1]. Instead, a characterization is studied of the set of computed instances of an atomic query, in a special case when the set is finite and the answers are ground. In the paper [13] of Kowalski completeness is discussed, but the example proofs concern only correctness. As a sufficient condition for completeness of a program $P$ he suggests $P \vdash T_S$, where $T_S$ is a specification in a form of
a logical theory. The condition seems impractical as it fails when $T_S$ contains auxiliary predicates, not occurring in $P$. It also requires that all the models of $P$ (including the Herbrand base) are models of the specification. But it seems that such specifications often have a substantially restricted class of models, maybe a single Herbrand model, cf. [4].

Deville [4] provides an approach where correctness and completeness of programs should follow from construction. No direct sufficient criteria for completeness, applicable to arbitrary programs, are given. Also the approach is not declarative, as it is based on an operational semantics of SLDNF-resolution.

Stärk [16] presents an elegant method of reasoning about a broad class of properties of programs with negation, executed under LDNF-resolutions. A tool to verify proofs mechanically was provided. The approach involves a rather complicated induction scheme, so it seems impossible to apply the method informally by programmers. Also, the approach is not fully declarative, as the order of literals in clause bodies is important.

A declarative sufficient condition for program completeness was given by DeRansart and Maluszynski [3]. The approach presented here stems from [8], the differences are discussed in the full version of this paper [7]. The main contribution from the former version [6] is proving completeness of pruned SLD-trees. The author is not aware of any other work on this issue.

**Preliminaries.** We use the standard notation and definitions [1]. An atom whose predicate symbol is $p$ will be called a $p$-atom (or an atom for $p$). Similarly, a clause whose head is a $p$-atom is a clause for $p$. In a program $P$, by procedure $p$ we mean the set of the clauses for $p$ in $P$.

We assume a fixed alphabet with an infinite set of function symbols. The Herbrand universe will be denoted by $\mathcal{HU}$, the Herbrand base by $\mathcal{HB}$, and the sets of all terms, respectively atoms, by $\mathcal{TU}$ and $\mathcal{TB}$. For an expression (a program) $E$ by $\text{ground}(E)$ we mean the set of ground instances of $E$ (ground instances of the clauses of $E$). $\mathcal{MP}$ denotes the least Herbrand model of a program $P$.

By “declarative” (property, reasoning, . . .) we mean referring only to logical reading of programs, thus abstracting from any operational semantics. In particular, properties depending on the order of atoms in clauses will not be considered declarative (as they treat equivalent conjunctions differently).

By a computed (respectively correct) answer for a program $P$ and a query $Q$ we mean an instance $Q\theta$ of $Q$ where $\theta$ is a computed (correct) answer substitution [1] for $Q$ and $P$. We often say just answer as each computed answer is a correct one, and each correct answer (for $Q$) is a computed answer (for $Q$ or for some its instance). Thus, by soundness and completeness of SLD-resolution, $Q\theta$ is an answer for $P$ iff $P \models Q\theta$.

Names of variables begin with an upper-case letter. We use the list notation of Prolog. So $[t_1, \ldots, t_n]$ ($n \geq 0$) stands for the list of elements $t_1, \ldots, t_n$. Only a term of this form is considered a list. (Thus terms like $[a, a|X]$, or $[a, a|a]$, where $a$ is a constant, are not lists). The set of natural numbers will be denoted by $\mathbb{N}$; $f: A \leftrightarrow B$ states that $f$ is a partial function from $A$ to $B$. 


The next section introduces the basic notions of specifications, correctness and completeness. Also, advantages of approximate specifications are discussed. After a brief overview of proving correctness, we discuss proving program completeness. Sect. 4 deals with proving that completeness is preserved under pruning. We finish with a discussion. For missing proofs, more examples etc see [7].

2 Correctness and completeness

Specifications. The purpose of a logic program is to compute a relation, or a few relations. A specification should describe these relations. It is convenient to assume that the relations are over the Herbrand universe. To describe such relations, one relation corresponding to each procedure of the program (i.e. to a predicate symbol), it is convenient to use a Herbrand interpretation. Thus a (formal) specification is a Herbrand interpretation, i.e. a subset of $\mathcal{HB}$.

Correctness and completeness. In imperative and functional programming, correctness usually means that the program results are as specified. In logic programming, due to its non-deterministic nature, we actually have two issues: correctness (all the results are compatible with the specification) and completeness (all the results required by the specification are produced). In other words, correctness means that the relations defined by the program are subsets of the specified ones, and completeness means inclusion in the opposite direction. In terms of specifications and the least Herbrand models we define:

**Definition 1** Let $P$ be a program and $S \subseteq \mathcal{HB}$ a specification. $P$ is correct w.r.t. $S$ when $M_P \subseteq S$; it is complete w.r.t. $S$ when $M_P \supseteq S$.

We will sometimes skip the specification when it is clear from the context. If a program $P$ is correct and complete w.r.t. a specification $S$ then, obviously, $M_P = S$. A program $P$ is correct w.r.t. a specification $S$ iff $Q$ being an answer of $P$ implies $S \models Q$. (Remember that $Q$ is an answer of $P$ iff $P \models Q$.) The program is complete w.r.t. $S$ iff $S \models Q$ implies that $Q$ is an answer of $P$. (Here our assumption on an infinite set of function symbols is needed [7].)

It is sometimes useful to consider local versions of these notions:

**Definition 2** A predicate $p$ in $P$ is correct w.r.t. $S$ when each $p$-atom of $M_P$ is in $S$, and complete w.r.t. $S$ when each $p$-atom of $S$ is in $M_P$.

An answer $Q$ is correct w.r.t. $S$ when $S \models Q$.

$P$ is complete for a query $Q$ w.r.t. $S$ when $S \models Q\theta$ implies that $Q\theta$ is an answer for $P$, for any ground instance $Q\theta$ of $Q$.

Informally, $P$ is complete for $Q$ when all the answers for $Q$ required by the specification $S$ are answers of $P$. Note that a program is complete w.r.t. $S$ iff it is complete w.r.t. $S$ for any query $Q \subseteq S$. 

Approximate specifications. Often it is difficult, and not necessary, to specify the relations defined by a program exactly; more formally, to require that $M_P$ is equal to a given specification. Often the relations defined by programs are not exactly those intended by programmers. For instance this concerns the programs in Chapter 3.2 of the textbook [17] defining predicates member/2, append/3, sublist/2, and some others. The defined relations are not those of list membership, concatenation, etc. However this is not an error, as for all intended queries the answers are as for a program defining the intended relations. The exact semantics of the programs is not explained in the textbook; such explanation is not needed. Let us look more closely at append/3.

Example 3. The program APPEND

\[
\{ \textit{app}([H|K], L, [H|M]) ← \textit{app}(K, L, M). \quad \textit{app}([], L, L). \} \]

does not define the relation of list concatenation. For instance, APPEND $\models \textit{app}([], 1, 1)$. In other words, APPEND is not correct w.r.t.

\[
S^0_{\text{APPEND}} = \{ \textit{app}(k, l, m) \in \textit{HB} \mid k, l, m \text{ are lists, } k \ast l = m \},
\]

where \(k \ast l\) stands for the concatenation of lists \(k, l\). It is however complete w.r.t. \(S^0_{\text{APPEND}}\), and correct w.r.t.

\[
S_{\text{APPEND}} = \{ \textit{app}(k, l, m) \in \textit{HB} \mid \text{if } l \text{ or } m \text{ is a list then } \textit{app}(k, l, m) \in S^0_{\text{APPEND}} \}.
\]

Correctness w.r.t. \(S_{\text{APPEND}}\) and completeness w.r.t. \(S^0_{\text{APPEND}}\) are sufficient to show that APPEND will produce the required results when used to concatenate or split lists. More precisely, the answers for a query \(Q = \textit{app}(s, t, u)\), where \(t\) is a list or \(u\) is a list, are \(\textit{app}(s\theta, t\theta, u\theta)\), where \(s\theta, t\theta, u\theta\) are lists and \(s\theta \ast t\theta = u\theta\). (The lists may be non-ground.)

Similarly, the procedures member/2 and sublist/2 are complete w.r.t. specifications describing the relation of list membership, and the sublist relation. It is easy to provide specifications, w.r.t. which the procedures are correct. For instance, member/2 is correct w.r.t. \(S_{\text{MEMBER}} = \{ \text{member}(t, u) \in \textit{HB} \mid \text{if } u = [t_1, \ldots, t_n] \text{ for some } n \geq 0 \text{ then } t = t_i, \text{ for some } 0 < i \leq n \}\).

The exact relations defined by programs are often misunderstood. For instance, in [5, Ex. 15] it is claimed that a program \(\text{Prog}_1\) defines the relation of list inclusion. In our terms, this means that predicate \(\text{included}\) of \(\text{Prog}_1\) is correct and complete w.r.t.

\[
\{ \text{included}(l_1, l_2) \in \textit{HB} \mid \text{l}_1, \text{l}_2 \text{ are lists, every element of } \text{l}_1 \text{ belongs to } \text{l}_2 \}.
\]

However the correctness does not hold: The program contains a unary clause \(\text{included}([], L)\), so \(\text{Prog}_1 \models \text{included}([], t)\) for any term \(t\).

The examples show that in many cases it is unnecessary to know the semantics of a program exactly. Instead it is sufficient to describe it approximately. An approximate specification is a pair of specifications \(S_{\text{compl}}, S_{\text{corr}}\).
for completeness and correctness. The intention is that the program is complete w.r.t. the former, and correct w.r.t. the latter: $S_{\text{compl}} \subseteq M_P \subseteq S_{\text{corr}}$. In other words, the specifications $S_{\text{compl}}, S_{\text{corr}}$ describe, respectively, which atoms have to be computed, and which are allowed to be computed. For the atoms from $S_{\text{corr}} \setminus S_{\text{compl}}$ the semantics of the program is irrelevant. By abuse of terminology, $S_{\text{corr}}$ or $S_{\text{compl}}$ will sometimes also be called approximate specifications.

**Proving correctness** We briefly discuss proving correctness, as it is complementary to the main subject of this paper. The approach is due to Clark [2].

**Theorem 4 (Correctness).** A sufficient condition for a program $P$ to be correct w.r.t. a specification $S$ is

for each ground instance $H \leftarrow B_1, \ldots, B_n$ of a clause of the program, if $B_1, \ldots, B_n \in S$ then $H \in S$.

**Example 5.** Consider a program SPLIT and a specification describing how the sizes of the last two arguments of $s$ are related ($|l|$ denotes the length of a list $l$):

$$s([],[],[]).$$ (1)

$$s([X|X]s,[X|Y]s,Zs) \leftarrow s(Xs,Zs,Ys).$$ (2)

$$S = \{ s(l,l_1,l_2) \mid l,l_1,l_2 \text{ are lists, } 0 \leq |l_1| - |l_2| \leq 1 \}.$$ (3)

SPLIT is correct w.r.t. $S$, by Th. 4 (the details are left for the reader, or see [7]). A stronger specification for which SPLIT is correct is shown in Ex. 11.

The sufficient condition is equivalent to $S \models P$, and to $T_P(S) \subseteq S$.

Notice that the proof method is declarative. The method should be well known, but is often neglected. For instance it is not mentioned in [1], where a more complicated method, moreover not declarative, is advocated. That method is not more powerful than the one of Th. 4 [8]. See [8] for further examples, explanations, references and discussion.

3 Proving completeness

We first introduce a notion of semi-completeness, and sufficient conditions under which semi-completeness of a program implies its completeness. Then a sufficient condition follows for semi-completeness. We conclude the section with a way of showing completeness directly without employing semi-completeness.

**Definition 6** A level mapping is a function $|\cdot| : \mathcal{HS} \to \mathbb{N}$ assigning natural numbers to atoms.

A program $P$ is **recurrent** w.r.t. a level mapping $|\cdot| \mid \mathcal{H} \mathcal{B}$ if, in every ground instance $H \leftarrow B_1, \ldots, B_n \in \text{ground}(P)$ of its clause ($n \geq 0$), $|H| > |B_i|$ for all $i = 1, \ldots, n$. A program is recurrent if it is recurrent w.r.t. some level mapping.

A program $P$ is **acceptable** w.r.t. a specification $S$ and a level mapping $|\cdot|$ if $P$ is correct w.r.t. $S$, and for every $H \leftarrow B_1, \ldots, B_n \in \text{ground}(P)$ we have $|H| > |B_i|$ whenever $S \models B_1, \ldots, B_{i-1}$. A program is acceptable if it is acceptable w.r.t. some level mapping and some specification.
The definition of acceptable is more general than that of \cite{1}, which requires \( S \) to be a model of \( P \). Both definitions make the same programs acceptable \cite{7}.

**Definition 7** A program \( P \) is **semi-complete** w.r.t. a specification \( S \) if \( P \) is complete w.r.t. \( S \) for any query \( Q \) for which there exists a finite SLD-tree.

In a less formal setting, the existence of a finite SLD-tree means that \( P \) with \( Q \) terminates under some selection rule. For a semi-complete program, if a computation for a query \( Q \) terminates then all the required by the specification answers for \( Q \) have been obtained. Note that a complete program is semi-complete. We also have:

**Proposition 8 (Completeness)** Let a program \( P \) be semi-complete w.r.t. \( S \). The program is complete w.r.t. \( S \) if

1. for each query \( A \in S \) there exists a finite SLD-tree, or
2. the program is recurrent, or
3. the program is acceptable (w.r.t. a specification \( S^* \) possibly distinct from \( S \)).

**Proving semi-completeness.** We need the following notion.

**Definition 9** A ground atom \( H \) is **covered by a clause** \( C \) w.r.t. a specification \( S \) \cite{13} if \( H \) is the head of a ground instance \( H ← B_1, \ldots, B_n \ (n \geq 0) \) of \( C \), such that all the atoms \( B_1, \ldots, B_n \) are in \( S \). A ground atom \( H \) is **covered by a program** \( P \) if it is covered by some clause \( C \in P \).

For instance, given a specification \( S = \{ p(s'(0)) \mid i \geq 0 \} \), atom \( p(s(0)) \) is covered both by \( p(s(X)) ← p(X) \) and by \( p(X) ← p(s(X)) \).

Now we present a sufficient condition for semi-completeness. Together with Prop.\cite{8} it provides a sufficient condition for completeness.

**Theorem 10 (Semi-completeness).** If all the atoms from a specification \( S \) are covered w.r.t. \( S \) by a program \( P \) then \( P \) is semi-complete w.r.t. \( S \).

**Example 11.** We show that program \( \text{SPLIT} \) from Ex.\cite{5} is complete w.r.t. \( S \).

\[
\text{SPLIT} = \left\{ \begin{array}{ll}
s(t_1, \ldots, t_{2n}, [t_1, \ldots, t_{2n-1}], [t_2, \ldots, t_{2n}]), & n \geq 0, \\
s([t_1, \ldots, t_{2n+1}], [t_1, \ldots, t_{2n}], [t_2, \ldots, t_{2n+1}]), & n \geq 0 \in \mathcal{HU} \end{array} \right\},
\]

where \([t_k, \ldots, t_l] \) denotes the list \([t_k, t_{k+2}, \ldots, t_l] \), for \( k, l \) both odd or both even.

Atom \( s([], [], []) \in \text{SPLIT} \) is covered by clause \cite{1}. For \( n > 0 \), any atom \( A = s([t_1, \ldots, t_{2n}], [t_1, \ldots, t_{2n-1}], [t_2, \ldots, t_{2n}]) \) is covered by an instance of \cite{2} with a body \( B = s([t_2, \ldots, t_{2n}], [t_2, \ldots, t_{2n}], [t_3, \ldots, t_{2n-1}]) \). Similarly, for \( n \geq 0 \) and any atom \( A = s([t_1, \ldots, t_{2n+1}], [t_1, \ldots, t_{2n+1}], [t_2, \ldots, t_{2n}]) \), the corresponding body is \( B = s([t_2, \ldots, t_{2n+2}], [t_2, \ldots, t_{2n}], [t_3, \ldots, t_{2n+2}]) \). In both cases, \( B \in \text{SPLIT} \).

(To see this, rename each \( t_i \) as \( t_{i+1} \).) So \( \text{SPLIT} \) is covered by \( \text{SPLIT} \). Thus \( \text{SPLIT} \) is semi-complete w.r.t. \( \text{SPLIT} \), by Th.\cite{10}.
Now by Prop. 9 the program is complete, as it is recurrent under the level mapping $|s(t, t_1, t_2)| = |t|$, where $|h| = 1 + |t|$ and $|f(t_1, \ldots, t_n)| = 0$ (for any ground terms $h, t, t_1, \ldots, t_n$, and any function symbol $f$ distinct from $\{ \mid \}$).

By Th. 12 the program is also correct w.r.t. $S_{\text{SPLIT}}$, as $S_{\text{SPLIT}} \models \text{SPLIT}$. (The details are left to the reader.) Hence $S_{\text{SPLIT}} = M_{\text{SPLIT}}$.

Note that the sufficient condition of Th. 10 is equivalent to $S \subseteq T_P(S)$, which implies $S \subseteq \text{gfp}(T_P)$.

The notion of semi-completeness is tailored for finite programs. An SLD-tree for a query $Q$ and an infinite program $P$ may be infinite, but with all branches finite. In such case, if the condition of Th. 10 holds then $P$ is complete for $Q$.

**Proving completeness directly.** Here we present another declarative way of proving completeness; a condition is added to Th. 10 so that completeness is implied directly. This also works for non-terminating programs. However when termination has to be shown anyway, applying Th. 10 seems simpler.

In this section we allow that a level mapping is a partial function $\mid H \mid : H \leftrightarrow \mathbb{N}$ assigning natural numbers to some atoms.

**Definition 12** A ground atom $H$ is recurrently covered by a program $P$ w.r.t. a specification $S$ and a level mapping $\mid \cdot \mid : H \leftrightarrow \mathbb{N}$ if $H$ is the head of a ground instance $H \leftarrow B_1, \ldots, B_n$ ($n \geq 0$) of a clause of the program, such that $|H|, |B_1|, \ldots, |B_n|$ are defined, $B_1, \ldots, B_n \in S$, and $|H| > |B_i|$ for all $i = 1, \ldots, n$.

For instance, given a specification $S = \{ p(s'(0)) \mid i \geq 0 \}$, atom $p(s(0))$ is recurrently covered by a program $\{ p(s(X)) \leftarrow p(X) \}$ under a level mapping for which $|p(s'(0))| = i$. No atom is recurrently covered by $\{ p(X) \leftarrow p(X) \}$. Obviously, if $H$ is recurrently covered by $P$ then it is covered by $P$. If $H$ is covered by $P$ w.r.t. $S$ and $P$ is recurrent w.r.t. $\mid \cdot \mid$ then $H$ is recurrently covered w.r.t. $S$, $\mid \cdot \mid$. The same holds for $P$ acceptable w.r.t. an $S' \supseteq S$.

**Theorem 13 (Completeness 2).** (A reformulation of Th. 6.1 of [3].) If, under some level mapping $\mid \cdot \mid : H \leftrightarrow \mathbb{N}$, all the atoms from a specification $S$ are recurrently covered by a program $P$ w.r.t. $S$ then $P$ is complete w.r.t. $S$.

**Example 14.** Consider a directed graph $E$. As a specification for a program describing reachability in $E$, take $S = S_p \cup S_e$, where

$$S_p = \{ p(t, u) \mid \text{there is a path from } t \text{ to } u \text{ in } E \},$$

$$S_e = \{ e(t, u) \mid (t, u) \text{ is an edge in } E \}.$$

Let $P$ consist of a procedure $p$: $p(X, X) \leftarrow p(X, Z) \leftarrow e(X, Y), p(Y, Z)$, and a procedure $e$ which is a set of unary clauses describing the edges of the graph. Assume the latter is complete w.r.t. $S_e$. Notice that when $E$ has cycles then infinite SLD-trees cannot be avoided, and completeness of $P$ cannot be shown by Prop. 9.
To apply Th.13 let us define a level mapping for the elements of $S$ such that $|e(t,u)| = 0$ and $|p(t,u)|$ is the length of a shortest path in $E$ from $t$ to $u$ (so $|p(t,t)| = 0$). Consider a $p(t,u) \in S$ where $t \neq u$. Let $t = t_0, t_1, \ldots, t_n = u$ be a shortest path from $t$ to $u$. Then $e(t,t_1), p(t_1,u) \in S$, $|p(t,u)| = n$, $|e(t,t_1)| = 0$, and $|p(t_1,u)| = n - 1$. Thus $p(t,u)$ is recurrently covered by $P$ w.r.t. $S$ and $| |$. The same trivially holds for the remaining atoms of $S$. So $P$ is complete w.r.t. $S$.

4 Pruning SLD-trees and completeness

Pruning some parts of SLD-trees is often used to improve efficiency of programs. It is implemented by using the cut, the if-then-else construct of Prolog, or built-ins, like once/1. Pruning preserves the correctness of a logic program, it also preserves termination under a given selection rule, but may violate the program’s completeness. We now discuss proving that completeness is preserved.

By a pruned SLD-tree for a program $P$ and a query $Q$ we mean a tree with the root $Q$ which is a connected subgraph of an SLD-tree for $P$ and $Q$. By an answer of a pruned SLD-tree we mean the computed answer of a successful SLD-derivation which is a branch of the tree. We will say that a pruned SLD-tree $T$ with root $Q$ is complete w.r.t. a specification $S$ if, for any ground $Q\theta$, $S \models Q\theta$ implies that $Q\theta$ is an instance of an answer of $T$. Informally, such a tree produces all the answers for $Q$ required by $S$.

We present two approaches for proving completeness of pruned SLD-trees. The first one is based on viewing pruning as skipping certain clauses while building the children of a node. The other deals with a restricted usage of the cut.

Pruning as clause selection. To facilitate reasoning about the answers of pruned SLD-trees, we will now view pruning as applying only certain clauses while constructing the children of a node. So we introduce subsets $\Pi_1, \ldots, \Pi_n$ of $P$. The intention is that $P = \bigcup_i \Pi_i$, and for each node the clauses of exactly one $\Pi_i$ are used. Programs $\Pi_1, \ldots, \Pi_n$ may be not disjoint.

Definition 15 Given programs $\Pi_1, \ldots, \Pi_n$ ($n > 0$), a c-selection rule is a function assigning to a query $Q'$ an atom $A$ in $Q'$ and one of the programs $\emptyset, \Pi_1, \ldots, \Pi_n$.

A csSLD-tree (cs for clause selection) for a query $Q$ and programs $\Pi_1, \ldots, \Pi_n$, via a c-selection rule $R$, is constructed as an SLD-tree, but for each node its children are constructed using the program selected by the c-selection rule. An answer of a csSLD-tree is defined in the expected way.

A c-selection rule may choose the empty program, thus making a given node a leaf. Notice that a csSLD-tree for $Q$ and $\Pi_1, \ldots, \Pi_n$ is a pruned SLD-tree for $Q$ and $\bigcup_i \Pi_i$. Conversely, for each pruned SLD-tree $T$ for $Q$ and a (finite) program $P$ there exist $n > 0, \Pi_1, \ldots, \Pi_n$ such that $P = \bigcup_i \Pi_i$ and $T$ is a csSLD-tree for $Q$ and $\Pi_1, \ldots, \Pi_n$. 
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Example 16. We show that completeness of each of $\Pi_1, \ldots, \Pi_n$ is not sufficient for completeness of a csSLD-tree for $\Pi_1, \ldots, \Pi_n$. Consider a program $P$:

$$q(X) \leftarrow p(Y, X). \quad (3) \quad p(a, s(X)) \leftarrow p(a, X). \quad (5)$$

and programs $\Pi_1 = \{3, 4, 6\}$, $\Pi_2 = \{3, 4, 5\}$, As a specification for completeness consider $S_0 = \{q(s^j(0)) \mid j \geq 0\}$. Each of the programs $\Pi_1, \Pi_2, P$ is complete w.r.t. $S_0$. Assume a c-selection rule $R$ choosing alternatively $\Pi_1, \Pi_2$ along each branch of a tree. Then the csSLD-tree for $q(s^j(0)) \in S_0$ via $R$ (where $j > 2$) has no answers, thus the tree is not complete w.r.t. $S_0$.

Consider a program $P = \bigcup_{i=1}^{n} \Pi_i$ and a specification $S = \bigcup_{i=1}^{n} S_i$. The intention is that each $S_i$ describes which answers are to be produced by using $\Pi_i$ in the first resolution step. We will sometimes call $\Pi_1, \ldots, \Pi_n, S_1, \ldots, S_n$ a split (of $P$ and $S$). Note that $\Pi_1, \ldots, \Pi_n$ or $S_1, \ldots, S_n$ may not be disjoint.

Definition 17 Let $S = \Pi_1, \ldots, \Pi_n, S_1, \ldots, S_n$ be a split, and $S = \bigcup S_i$.

Specification $S_i$ is suitable for an atom $A$ w.r.t. $S$ when no instance of $A$ is in $S \setminus S_i$. (In other words, when ground$(A) \cap S \subseteq S_i$.) We also say that a program $\Pi_i$ is suitable for $A$ w.r.t. $S$ when $S_i$ is.

A c-selection rule is compatible with $S$ if for each non-empty query $Q$ it selects an atom $A$ and a program $\Pi_i$ such that

- $\Pi_i \in \{\Pi_1, \ldots, \Pi_n\}$ is suitable for $A$ w.r.t. $S$, or
- none of $\Pi_1, \ldots, \Pi_n$ is suitable for $A$ w.r.t. $S$ and $\Pi = \emptyset$ (so $Q$ is a leaf).

A csSLD-tree for $\Pi_1, \ldots, \Pi_n$ via a c-selection rule compatible with $S$ is said to be weakly compatible with $S$. The tree is compatible with $S$ if a nonempty program is selected for each of its nonempty node.

The intuition is that when $\Pi_i$ is suitable for $A$ then $S_i$ is a fragment of $S$ sufficient to deal with $A$. It describes all the answers for query $A$ required by $S$.

The reason of incompleteness of the trees in Ex. 16 may be understood as selecting a $\Pi_i$ not suitable for the selected atom. Take $S = \Pi_1, \Pi_2, S_0 \cup S_1, S_0 \cup S_2$, where $S_i = \{p(b, s^i(0)) \mid i \geq 0\}$ and $S_0 = \{p(a, s^i(0)) \mid i \geq 0\}$. In the incomplete trees, $\Pi_1$ is selected for an atom $A = p(a, u)$, or $\Pi_2$ is selected for an atom $B = p(b, u)$ (where $u \in TU$). However $\Pi_1$ is not suitable for $A$ whenever $A$ has an instance in $S$ (as then ground$(A) \cap S \subseteq S_0 \cup S_1$); similarly for $\Pi_2$ and $B$.

When $\Pi_i$ is suitable for $A$ then if each atom of $S_i$ is covered by $\Pi_i$ (w.r.t. $S$), then using for $A$ only the clauses of $\Pi_i$ does not impair completeness w.r.t. $S$.

Theorem 18. Let $P = \bigcup_{i=1}^{n} \Pi_i$ (where $n > 0$) be a program, $S = \bigcup_{i=1}^{n} S_i$ a specification, and $T$ a csSLD-tree for $\Pi_1, \ldots, \Pi_n$. If

1. for $i = 1, \ldots, n$, all the atoms from $S_i$ are covered by $\Pi_i$ w.r.t. $S$, and
2. $T$ is compatible with $\Pi_1, \ldots, \Pi_n, S_1, \ldots, S_n$,
3. (a) $T$ is finite, or
   - (b) program $P$ is recurrent,
(c) $P$ is acceptable (possibly w.r.t. a specification distinct from $S$), and $T$ is built under the Prolog selection rule 

then $T$ is complete w.r.t. $S$.

**Example 19.** The following program SAT0 is a simplification of a fragment of the SAT solver of [11] discussed in [6]. Pruning is crucial for the efficiency and usability of the original program.

\[ p(P-P, []) \quad (7) \quad q(V-P, \_ ) \leftarrow V = P. \quad (10) \]
\[ p(V-P, [B|T]) \leftarrow q(V-P, [B|T]). \quad (8) \]
\[ q(\_, [A|T]) \leftarrow p(A, T). \quad (11) \]
\[ p(V-P, [B|T]) \leftarrow q(B, [V-P|T]). \quad (9) \]
\[ P = P. \quad (12) \]

The program is complete w.r.t. a specification

\[ S = \{ p(t_0-u_0, [t_1-u_1, \ldots, t_n-u_n]), \quad n \geq 0, \ t_0, \ldots, t_n, u_0, \ldots, u_n \in T, \ t_i = u_i \text{ for some } i \in \{0, \ldots, n\} \} \cup S_u \]

where $T = \{ \text{false, true} \} \in HUA$, and $S_u = \{ t = t \mid t \in HUA \}$. We omit a completeness proof, mentioning only that SAT0 is recurrent w.r.t. a level mapping $|p(t, u)| = 2|u| + 2$, $|q(t, u)| = 2|u| + 1$, $|=t, u| = 0$, where $|u|$ is as in Ex. [11].

The first case of pruning is due to redundancy within (8), (9); both $\Pi_1 = SAT0 \setminus \{9\}$ and $\Pi_2 = SAT0 \setminus \{8\}$ are complete w.r.t. $S$. For any selected atom at most one of (8), (9) is to be used, and the choice is dynamic. As the following reasoning is independent from this choice, we omit further explanations.

So in such pruned SLD-trees the children of each node are constructed using one of programs $\Pi_1, \Pi_2$. Thus they are csSLD-trees for $\Pi_1, \Pi_2$. They are compatible with $S = \Pi_1, \Pi_2, S, S$ (as $\Pi_1, \Pi_2$ are trivially suitable for any $A$, due to $S_t = S$ and $S \setminus S_t = \emptyset$ in Def. [17]). Each atom of $S$ is covered w.r.t. $S$ both by $\Pi_1$ and $\Pi_2$. As SAT0 is recurrent, by Th. [18] each such tree is complete w.r.t. $S$.

**Example 20.** We continue with program SAT0 and specification $S$ from the previous example, and add a second case of pruning. When the selected atom is of the form $A = q(s_1, s_2)$ with a ground $s_1$ then only one of clauses (10), (11) is needed $-$ (10) when $s_1$ is of the form $t-t$, and (11) otherwise. The other clause can be abandoned without losing the completeness w.r.t. $S$.

Actually, SAT0 is included in a bigger program, say $P = SAT0 \cup S_0$. We skip the details of $S_0$, let us only state that $P$ is recurrent, $S_0$ does not contain any clause for $p$ or for $q$, and that $P$ is complete w.r.t. a specification $S' = S \cup S_0$ where $S_0$ does not contain any $p$- or $q$-atom. (Hence each atom of $S_0$ is covered by $H_0$ w.r.t. $S'$.)

---

1 The same holds for $A$ of the form $q(s_{11}, s_{12})$, or $q(s_{11}, s_{12}, s_2)$ with non-unifiable $s_{11}$, $s_{12}$. The pruning is implemented using the if-then-else construct in Prolog: $q(V-P, [A|T]) ::= \forall P \rightarrow \text{true} ; p(A, T)$. (And the first case of pruning by $p(V-P, [B|T]) ::= \text{nonvar}(V) \rightarrow q(V-P, [B|T])$ ; $q(B, [V-P|T])$.)
To formally describe the trees for $P$ resulting from both cases of pruning, consider $S = I_0, \ldots, I_5, S_0, \ldots, S_5$, where

\begin{align*}
I_1 &= \{(7), (8)\}, & I_2 &= \{(7), (9)\}, & S_1 &= S_2 = S \cap \{p(s, u) \mid s, u \in HU\}, \\
I_3 &= \{(10)\}, & S_3 &= S \cap \{q(t, s) \mid t, s \in HU\}, \\
I_4 &= \{(11)\}, & S_4 &= S \cap \{q(t, u, s) \mid t, u, s \in HU, t \neq u\}, \\
I_5 &= \{(12)\}, & S_5 &= S. \\
\end{align*}

Each atom from $S_i$ is covered by $I_i$ w.r.t. $S$ (for $i = 0, \ldots, 5$). For each $q$-atom with its first argument ground, $I_3$ or $I_4$ (or both) is suitable. For each remaining atom from $T\cup B$, a program from $I_0, I_1, I_2, I_5$ is suitable.

Consider a pruned SLD-tree $T$ for $P$ (employing the two cases of pruning described above). Assume that each $q$-atom selected in $T$ has its first argument ground. Then $T$ is a csSLD-tree compatible with $S$. From Th. 18 it follows that $T$ is complete w.r.t. $S$.

The restriction on the selected $q$-atoms can be implemented by means of Prolog delays. This can be done in such a way that floundering is avoided for the intended initial queries. So the obtained pruned trees are as $T$ above, and the pruning preserves completeness of the program.

**Example 21.** A Prolog program \{\texttt{nop(adam, 0):-!}. \texttt{nop(eve, 0):-!}. \texttt{nop(X, 2)}.\} is an example of difficulties and dangers of using the cut in Prolog. Due to the cut, for an atomic query $A$ only the first clause with the head unifiable with $A$ will be used. The program can be seen as logic program $P = I_1 \cup I_2 \cup I_3$ executed with pruning, where (for $i = 1, 2, 3$) $I_i$ is the $i$-th clause of the program with the cut removed. The intended meaning is $S = S_1 \cup S_2 \cup S_3$, where $S_1 = \{\texttt{nop(adam, 0)}\}$, $S_2 = \{\texttt{nop(eve, 0)}\}$, and $S_3 = \{\texttt{nop(t, 2) \in HB} \mid t \notin \{\texttt{adam, eve}\}\}$. Note that all the atoms from $S_i$ are covered by $I_i$ (for $i = 1, 2, 3$). (We do not discuss here the (in)correctness of the program.)

Let $S$ be $I_1, I_2, I_3, S_1, S_2, S_3$. Consider a query $A = \texttt{nop(t, Y)}$ with a ground $t$. If $t = \texttt{adam}$ then only $I_1$ is suitable for $A$ w.r.t. $S$, if $t = \texttt{eve}$ then only $I_2$ is. For $t \notin \{\texttt{adam, eve}\}$ the suitable program is $I_3$. So for a query $A$ the pruning due to the cuts in the program results in selecting a suitable $I_i$, and the obtained csSLD-tree compatible with $S$. By Th. 18 the tree is complete w.r.t. $S$.

For a query $\texttt{nop(X, Y)}$ or $\texttt{nop(X, 0)}$ only the first clause, i.e. $I_1$, is used. However $I_1$ is not suitable for the query (w.r.t. $S$), and the csSLD-tree is not compatible with $S$. The tree is not complete (w.r.t. $S$).

**The cut in the last clause.** The previous approach is based on a somehow abstract semantics in which pruning is viewed as clause selection. Now we present an approach referring directly to Prolog with the cut. However the usage of the cut is restricted to the last clause of a procedure. Without this restriction, the cut can be used to implement negation as failure. So the author expects that completeness for such programs should be studied together with that for programs with negation. We consider LD-resolution (as interaction of the cut with delays introduces additional complications).
Example 24. Consider a program IN and specifications:

\[ \text{T} \]

then

\[ H_\rho \]

condition 3 holds for an substitution \( \theta \) to the cut occurring in the last clause(s) of some procedure(s) of \( \text{P} \). Let \( S \) be a specification, where \( S \) is an instance \([7]\). Note that if \( \theta \) is a ground atom \( \text{A} \) w.r.t. \( \text{pre,post} \). Let Theorem 23. Informally, condition 3 says that \( \text{A} \) could be produced out of each “related” answer for \( A_1, \ldots, A_{k-1}, \text{A}_k, \ldots, \text{A}_n \), and some answers for \( A_k, \ldots, A_n \), specified by \( S \). If condition 3 holds for an \( H_\rho \) then it holds for any its instance \( H_\rho \delta \) of which \( A \) is an instance \([7]\). Note that if \( A \) is adjustably covered by \( C \) w.r.t. \( S \), \( \text{pre,post} \) if \( A \) is covered by \( C \) and the cut does not occur in \( C \), or the following three conditions hold:

1. \( C = H \leftarrow A_1, \ldots, A_{k-1}, !, A_k, \ldots, A_n \),
2. \( A \) is covered by \( H \leftarrow A_1, \ldots, A_{k-1} \) w.r.t. \( S \),
3. – for any instance \( H_\rho \in \text{pre} \) such that \( A \) is an instance of \( H_\rho \),
   – for any ground instance \((A_1, \ldots, A_{k-1})_\rho \) such that \( A_1_\rho, \ldots, A_{k-1}_\rho \in \text{post} \),
   – there exists a ground instance \( C_\rho \sigma \) of \( C \) such that \( H_\rho \sigma = A \), and \( A_k_\rho \sigma, \ldots, A_n \rho \sigma \in S \),
where \( \text{dom}(\rho) \subseteq \text{vars}(H) \), \( \text{rng}(\rho) \cap \text{vars}(C) \subseteq \text{vars}(H) \), \( \text{dom}(\rho) \cap \text{rng}(\rho) = \emptyset \), and \( \text{dom}(\eta) = \text{vars}((A_1, \ldots, A_{k-1})_\rho) \).

We need to reason about the atoms selected in the derivations. So we employ a (non-declarative) approach to reason about LD-derivations, presented in \([1]\). A specification in this approach, let us call it call-success specification, is a pair \( \text{pre,post} \in \text{TB} \) of sets of atoms, closed under substitution. When a program is well-asserted \([1]\) w.r.t. \( \text{pre,post} \) then in each LD-derivation every selected atom is from \( \text{pre} \) and each corresponding computed answer is in \( \text{post} \), provided that the derivation begins with an atomic query from \( \text{pre} \). See \([1]\) for further explanations.

By \( \text{vars}(E) \) we denote the set of variables occurring in an expression \( E \). For a substitution \( \theta = \{X_1/t_1, \ldots, X_n/t_n\} \), \( \text{dom}(\theta) = \{X_1, \ldots, X_n\} \), and \( \text{rng}(\theta) = \text{vars}(\{t_1, \ldots, t_n\}) \).

Definition 22. Let \( S \) be a specification, and \( \text{pre,post} \) a call-success specification. A ground atom \( A \) is adjustably covered by a clause \( C \) w.r.t. \( S \) and \( \text{pre,post} \) if \( A \) is covered by \( C \) and the cut does not occur in \( C \), or the following three conditions hold:

1. \( C = H \leftarrow A_1, \ldots, A_{k-1}, !, A_k, \ldots, A_n \),
2. \( A \) is covered by \( H \leftarrow A_1, \ldots, A_{k-1} \) w.r.t. \( S \),
3. – for any instance \( H_\rho \in \text{pre} \) such that \( A \) is an instance of \( H_\rho \),
   – for any ground instance \((A_1, \ldots, A_{k-1})_\rho \) such that \( A_1_\rho, \ldots, A_{k-1}_\rho \in \text{post} \),
   – there exists a ground instance \( C_\rho \sigma \) of \( C \) such that \( H_\rho \sigma = A \), and \( A_k_\rho \sigma, \ldots, A_n \rho \sigma \in S \),
where \( \text{dom}(\rho) \subseteq \text{vars}(H) \), \( \text{rng}(\rho) \cap \text{vars}(C) \subseteq \text{vars}(H) \), \( \text{dom}(\rho) \cap \text{rng}(\rho) = \emptyset \), and \( \text{dom}(\eta) = \text{vars}((A_1, \ldots, A_{k-1})_\rho) \).

Informally, condition 3 says that \( A \) could be produced out of each “related” answer for \( A_1, \ldots, A_{k-1}, \text{A}_k, \ldots, \text{A}_n \), and some answers for \( A_k, \ldots, A_n \), specified by \( S \). If condition 3 holds for an \( H_\rho \) then it holds for any its instance \( H_\rho \delta \) of which \( A \) is an instance \([7]\). Note that if \( A \) is adjustably covered by \( C \) w.r.t. \( S \), \( \text{pre,post} \) if \( A \) is covered by \( C \) w.r.t. \( S \).

Theorem 23. Let \( S \) be an approximate specification, \( \text{pre,post} \) a call-success specification, where \( S \subseteq \text{post} \), and a program \( P \) be well-asserted w.r.t. \( \text{pre,post} \). Let \( T \) be a pruned LD-tree for \( P \) and an atomic query \( Q \), where pruning is due to the cut occurring in the last clause(s) of some procedure(s) of \( P \). If

– \( T \) is finite, \( Q \in \text{pre} \), and
– each \( A \in S \) is adjustably covered by a clause of \( P \) w.r.t. \( S \) and \( \text{pre,post} \)
then \( T \) is complete w.r.t. \( S \).

Example 24. Consider a program IN and specifications:

\[ \text{in}([], L), \quad \text{m}([E], [E|L]), \]
\[ \text{in}([H][T], L) \leftarrow \text{m}([H,L], !, \text{in}(T,L)), \quad \text{m}([E], [E|L]) \leftarrow \text{m}(E,L). \]

\[ S = \text{Sm} \cup \text{Sin}, \quad \text{pre} = \text{pre}_m \cup \text{pre}_in, \quad \text{post} = \text{post}_m \cup \text{post}_in, \quad \text{where} \]
\[ \text{pre}_m = \{ m(u,t) \in \text{TB} \mid t \text{ is a list} \}, \]
\[ \text{pre}_in = \{ \text{in}(u,t) \in \text{HB} \mid u,t \text{ are ground lists} \}, \]
\[ \text{post}_m = \{ m(t_1, [t_1, \ldots, t_n]) \in \text{TB} \mid 1 \leq i \leq n \}, \]
\[ \text{post}_in = \{ \text{in}([u_1, \ldots, u_m], [t_1, \ldots, t_n]) \in \text{HB} \mid \{u_1, \ldots, u_m\} \subseteq \{t_1, \ldots, t_n\} \}, \]
\[ \text{Sm} = \text{post}_m \cap \text{HB}, \quad \text{Sin} = \text{post}_in. \]
he program is well-asserted (we skip the details). We show that each atom \( A = \text{in}(u, t) \in S_{in} \), where \( u = [u_1, \ldots, u_m] \), \( t = [t_1, \ldots, t_n] \), \( m > 0 \), is adjustably covered by the second clause \( C \) of IN. Let \( C_0 \) be \( \text{in}([H|T], L) \leftarrow m(H, L) \). Now \( A \) is covered by \( C_0 \) w.r.t. \( S \) (\( A \leftarrow m(u_1, t) \) is a relevant ground instance of \( C_0 \)).

Any instance \( \text{in}([H|T], L) \rho \in \text{pre} \) of the head of \( C_0 \) is ground. Moreover the whole \( C_0 \rho \) is ground. So substitutions \( \eta, \sigma \) from Df. 22 are empty. If \( A \) is an instance of (thus equal to) \( \text{in}([H|T], L) \rho \) then \( \text{in}(T, L) \rho = \text{in}([u_2, \ldots, u_m], t) \in S. \) So condition 3 of Df. 22 is satisfied.

Thus \( A \) is adjustably covered by \( C \). It is easy to check that all the remaining atoms of \( S \) are covered by IN w.r.t. \( S \), and that IN is recurrent (for \( m(s, t) = |t| \), \( |\text{in}(s, t)| = |s| + |t| \), \( |t| \) as in Ex. 11). Thus each LD-tree for IN and a query \( Q \in \text{pre} \) is finite. By Th. 23, each such tree pruned due to the cut is complete w.r.t. \( S \). Notice that condition 3 does not hold when non-ground arguments of \( \text{in} \) are allowed in \( \text{pre} \), and that for such queries some answers may be pruned.

### 5 Discussion

**Declarativeness.** Without declarative ways of reasoning about correctness and completeness of programs, logic programming would not deserve to be called a declarative programming paradigm. The sufficient condition for proving correctness (Th. 2), that for semi-completeness of Th. 10, and those for completeness of Prop. 8.2 and Th. 13 are declarative. However the sufficient conditions for completeness of Prop. 8.1 and 8.3 are not, as they refer to program termination, or depend on the order of atoms in clause bodies.

Declarative completeness proofs employing Prop. 8.2 or Th. 13 imply termination, or require reasoning similar to that in termination proofs. So proving completeness by means of semi-completeness and termination may be a reasonable compromise between declarative and non-declarative reasoning, as termination has to be shown anyway in most of practical cases.

**Granularity of proofs.** Note that the sufficient condition for correctness deals with single clauses, that for semi-completeness – with procedures, and those for completeness take into account a whole program.

**Incompleteness diagnosis.** There is a close relation between completeness proving and incompleteness diagnosis 15. As the reason of incompleteness, a diagnosis algorithm finds an atom from \( S \) that is not covered by the program. Thus it finds a reason for violating the sufficient conditions for semi-completeness and completeness of Th. 10, 13.

**Approximate specifications.** We found that approximate specifications are crucial in avoiding unnecessary complications in dealing with correctness and completeness of programs (cf. Sect. 2, 8, 16, 17). For instance, in the main example of 6 (and in its simpler version in Ex. 19, 20) finding an exact specification is not easy, and is unnecessary. The required property of the program is described
more conveniently by an approximate specification. Moreover, as this example shows, in program development the semantics of (common predicates in) the consecutive versions of a program may differ. What is unchanged is correctness and completeness w.r.t. an approximate specification.

**Approximate specifications in program development.** This suggests a generalization of the paradigm of program development by semantics preserving program transformations [14]: it is useful and natural to use transformations which preserve correctness and completeness w.r.t. an approximate specification.

**Approximate specifications in debugging.** In declarative diagnosis [15] the programmer is required to know the exact intended semantics of the program. This is a substantial obstacle to using declarative diagnosis in practice. Instead, an approximate specification can be used, with the specification for correctness (respectively completeness) applied in incorrectness (incompleteness) diagnosis. See [7] for discussion and references.

**Interpretations as specifications.** This work uses specifications which are interpretations. (The same kind of specifications is used, among others, in [1], and in declarative diagnosis.) There are however properties which cannot be expressed by such specifications [8]. For instance one cannot express that some instance of an atomic query \( A \) should be an answer; one has to specify the actual instance(s).

Other approach is needed for such properties, possibly with specifications which are logical theories (where axioms like \( \exists X. A \) can be used).

**Applications.** We want to stress the simplicity and naturalness of the sufficient conditions for correctness (Th. 4) and semi-completeness (Th. 10, the condition is a part of each discussed sufficient condition for completeness). Informally, the first one says that the clauses of a program should produce only correct conclusions, given correct premises. The other says that each ground atom that should be produced by \( P \) has to be produced by a clause of \( P \) out of atoms produced by \( P \). The author believes that this is a way a competent programmer reasons about (the declarative semantics of) a logic program.

Paper [6] illustrates practical applicability of the methods presented here. It shows a systematic construction of a non-trivial Prolog program (the SAT solver of [11]). Starting from a formal specification, a definite clause logic program is constructed hand in hand with proofs of its correctness, completeness, and termination under any selection rule. The final Prolog program is obtained by adding control to the logic program (delays and pruning SLD-trees). Adding control preserves correctness and termination. However completeness may be violated by pruning, and by floundering related to delays. By Th. [18] the program with pruning remains complete. Proving non-floundering is outside of the scope of this work. See [9] for a related analysis algorithm, applicable in this case [12].

The example shows how well “logic” could be separated from “control.” The whole reasoning related to correctness and completeness can be done declaratively, abstracting from any operational semantics.

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2 In [6] a weaker version of Th. [18] was used, and one case of pruning was discussed informally. A proof covering both cases of pruning is illustrated here in Ex. [20].
Future work. A natural continuation is developing completeness proof methods for programs with negation (a first step was made in [8]), maybe also for constraint logic programming and CHR (constraint handling rules). Further examples of proofs for various programs are necessary. An interesting task is formalizing and automatizing the proofs, a first step is formalization of specifications.

Conclusion. Reasoning about completeness of logic program has been, surprisingly, almost neglected. This paper presents a few sufficient conditions for completeness. As an intermediate step we introduced a notion of semi-completeness. We also propose two methods of proving completeness in presence of pruning. The presented methods are, to a large extent, declarative. Examples suggest that the approach is applicable – maybe at informal level – in practice of Prolog programming.

References
Polynomial Approximation to Well-Founded Semantics for Logic Programs

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Abstract. The well-founded semantics of normal logic programs has two main utilities, one being an efficiently computable semantics with a unique intended model, and the other serving as polynomial time constraint propagation for the computation of answer sets of the same program. When logic programs are generalized to support constraints of various kinds, the semantics is no longer tractable, which makes the second utility doubtful. This paper considers the possibility of tractable but incomplete methods, which in general may miss information in the computed result, but never generates wrong conclusions. For this goal, we first formulate a well-founded semantics for logic programs with generalized atoms, which generalizes logic programs with arbitrary aggregates/constraints/dl-atoms. As a case study, we show that the method of removing non-monotone dl-atoms for the well-founded semantics by Eiter et al. actually falls into this category. We also present a case study for logic programs with standard aggregates.

Keywords: Polynomial Approximation, Well-Founded Semantics, Generalized Atoms.

1 Introduction

Logic programming with negation is a programming paradigm for declarative knowledge representation and problem solving [6]. In a logic program, a problem is represented as a set of rules and the solutions of the problem are determined under a predefined semantics. In more recent years, logic programs have been incorporated with various types of predefined atoms, to enable reasoning with constraints and external knowledge sources. Examples include logic program with aggregates [1, 13, 16, 23], logic programs with constraints, in an abstract form [17, 24] or in terms of global constraints [19], dl-programs [9, 11] and logic programs with external sources (e.g. HEX-programs [10]).

For a unifying framework for the study of logic programs, following [2], a constraint atom in its generic form is called a generalized atom and logic programs with these atoms are called logic programs with generalized atoms. The main interest of this paper lies in the two case studies of polynomial approximation to the well-founded semantics of these programs.
It is well-known that the well-founded semantics of normal logic programs [26] can be computed in polynomial time. Besides many applications, the mechanism of computing the well-founded semantics under a partial truth value assignment can be employed in the computation of answer sets. This is utilized, for example, in SMODELS [22] by the Expand function (though the term of well-founded semantics was not used explicitly at that time); it is invoked before and after every choice point in search. To serve as constraint propagation, it is essential that this computational process is effective. However, when constraints are incorporated into logic programs, the resulting well-founded semantics becomes intractable, since determining the satisfiability of a generalized atom by all consistent extensions of a partial interpretation is co-NP-complete [1, 23].

For the utility in constraint propagation in answer set computation, this is a problem.

In this paper, we consider tractable but incomplete methods. That is, we would like to transform a given program $P$ to another one $P'$ in polynomial time, so that

- the well-founded semantics of $P'$ is tractable, and
- the well-founded semantics of $P'$ specifies a subset of conclusions by the well-founded semantics of $P$ (when restricted to the language of $P$).

Of course, such a transformation should be non-trivial, as the empty program trivially satisfies these properties.

The approach described above is based on a single transformed program $P'$. In general, however, the computation of an approximation of the well-founded semantics may be carried out by a collection of program components each of which is employed for some specific computations, as long as the overall process takes polynomial time. We will call all of such approaches incomplete methods. But for simplicity, in the following we feel free to describe the effect of an incomplete method by referring to the approach based on a single transformed program $P'$.

Incomplete methods have practical implications. For example, for computing the well-founded semantics of a logic program with (arbitrary) generalized atoms, one can always compute such an approximation first. For the utility in answer set computation, let us assume that the well-founded semantics of $P$ approximates the answer sets of $P$, under a suitable definition of answer set. That is, any atom true (resp. false) in the well-founded semantics of $P$, called well-founded (resp. unfounded), remains to hold in any answer set of $P$. Then, in the computation of answer sets of $P$, constraint propagation can be performed by computing the well-founded semantics of $P'$ that extends a partial truth value assignment, at any choice point in search for answer set. The above conditions guarantee that constraint propagation is correct and effective.

It is well-known that an efficient constraint propagation mechanism is essential in all popular search engines, e.g., BCP in Boolean Satisfiability (SAT) [7], the Expand function in Answer Set Programming (ASP) [22], and various consistency techniques in Constraint Programming (CP) [21].

To study incomplete methods, the first question is which well-founded semantics is to be approximated, and how it is defined. There seems to be a consensus in the field about what such a well-founded semantics should be, at least for some specific classes.\footnote{This assumes that the satisfaction of a generalized atom by one interpretation can be determined in polynomial time.}
of logic programs. In [12] for (disjunctive) logic programs with arbitrary aggregates, the notion of unfounded set is defined but the well-founded semantics itself is not spelled out explicitly. A similar approach is presented in [11], where a well-founded semantics for programs with monotone and anti-monotone dl-atoms is defined, which is also based on the notion of unfounded set, followed by a comment (cf. Section 9.2) that the same approach can be generalized to programs with arbitrary dl-atoms. But the details are not provided. In [1], a study of logic programs with monotone and anti-monotone aggregates is presented. Despite these efforts, to our knowledge, there has been no unified approach such that all of the above classes of logic programs, the well-founded semantics falls into the same theoretical framework. Apparently, a unified approach is important in studying the common properties of logic programs in syntactically different forms. Thus, our first task is to formally define this well-founded semantics for logic programs with generalized atoms. In fact, it is in [2] that introduces the notion of generalized atom for the study of answer set semantics for various classes of logic programs. Hence, our work presented here can be viewed as complementing that of [2] for the study of well-founded semantics.

The well-founded semantics for logic programs with arbitrary generalized atoms is not tractable in the general case. We then carry out two case studies on the possibility of polynomial approximation. In the first, we consider dl-programs with arbitrary dl-atoms. In fact, the authors of [11] give a simple rewrite scheme for removing non-monotone dl-atoms from a program, in which every occurrence of the constraint operator $\cap$ is replaced. It is shown that the data complexity for $\cap$-free dl-programs is tractable. However, we show that the transformation is not faithful - the well-founded semantics of the transformed program in general is not equivalent to the well-founded semantics of the original program. We show however that the rewrite scheme is correct. Namely, though it may miss some conclusions, it never generates incorrect ones. This thus provides the first case of polynomial approximation for a practical class of logic programs with generalized atoms.

We then turn our attention to logic programs with aggregates. We adopt the disjunctive rewrite scheme of c-atoms described in [15] for aggregates with standard aggregate functions [23]. We show disjunctive rewrite that the disjunctive rewrite scheme is actually an incomplete method - it may miss conclusions but never generates wrong ones for logic programs with (standard) aggregates. We also show how to apply two different methods of disjunctive rewrite to construct a polynomial process for the computation of an approximation of the well-founded semantics, for logic programs with standard aggregates.

The rest of the paper is organized as follows: Section 2 provides syntax and notations. Section 3 defines the well-founded semantics for logic program with generalized atoms and studies its basic properties. Then, we carry out case studies on polynomial approximation for dl-programs in Section 4, and for logic programs with aggregates in Section 5, followed by related work and discussion in Section 6.
2 Preliminaries

2.1 Language

We assume a language $\mathcal{L}$ that includes a countable set of ground atoms $\Sigma$ (atoms for short). A literal is either an atom or its negation. An interpretation is a subset of $\Sigma$.

Following the spirit of [2], a generalized atom $\alpha$ on $\Sigma$ is a mapping from interpretations to the Boolean truth values $t$ and $f$.

There is no particular syntax requirement, except that a generalized atom $\alpha$ is associated with a domain, denoted $\text{Dom}(\alpha) \subseteq \Sigma$, and the size of $\alpha$ is $|\text{Dom}(\alpha)|$.

We assume that a generalized atom may only mention atoms in $\Sigma$ (i.e., it is not nested). An interpretation $I$ satisfies $\alpha$, written $I \models \alpha$, if $\alpha$ maps $I$ to $t$, $I$ does not satisfy $\alpha$, written $I \not\models \alpha$, if $\alpha$ maps $I$ to $f$.

Intuitively, a generalized atom is a constraint, whose semantics is defined externally by how it may be satisfied by sets of atoms.

Since an atom in $\Sigma$ can be viewed as a special case of a generalized atom, in the sequel, the term generalized atoms also includes all atoms in $\Sigma$. At times of possible confusion, we may use the term ordinary atoms for those in $\Sigma$, for distinction. As usual, an interpretation $I$ satisfies an ordinary atom $a$, written $I \models a$, if $a \in I$ ($\alpha$ maps $I$ to $t$), and $I$ does not satisfy $a$, written $I \not\models a$ (or $I \models \neg a$), if $a \not\in I$ ($\alpha$ maps $I$ to $f$).

A generalized atom $\alpha$ is monotone if for every interpretation $I$ such that $I \models \alpha$, we have $J \models \alpha$ for all interpretations $J$ such that $J \supseteq I$; it is anti-monotone if for every interpretation $I$ such that $I \models \alpha$, we have $J \not\models \alpha$ for all interpretations $J$ such that $J \subseteq I$; it is non-monotone if it is not monotone.

Aggregates atoms, abstract constraint atoms, HEX-atoms, dl-atoms, weight constraints, and global constraints in constraint satisfaction can all be modeled by generalized atoms. Below we sketch an example.

Example 1. – An aggregate atom consists of a mapping from multi-sets to numbers, a comparison operator, and a value. For example, following the notation of [23], $\text{SUM}([X]|p(X)) \geq 1$ denotes a (non-ground) aggregate atom: after grounding one can view the set of ground instances of $p(X)$ as the domain of the atom. The semantics of the aggregate is defined by interpretations in which the sum of the arguments in satisfied $p(.)$ is greater than or equal to 1.

An abstract constraint atom (or a c-atom) is of the form $(D, S)$, where $D \subseteq \Sigma$ serves as the domain and $S$ is a set of subsets of $D$ representing admissible solutions of the c-atom [18]. A c-atom $(D, S)$ is satisfied by an interpretation $I$ iff $I \cap D \in S$. For example, the aggregate atom $\text{SUM}([X]|p(X)) \geq 1$ above with domain $D = \{p(1), p(-1), p(2)\}$ can be represented by a c-atom $(D, S)$ where $S = \{\{p(1)\}, \{p(2)\}, \{p(-1), p(2)\}, \{p(1), p(2)\}, \{p(1), p(-1), p(2)\}\}$. But notice the difference between c-atoms and generalized atoms: a c-atom contains a specification of how it may be satisfied by sets of atoms internally, while the satisfiability of a generalized atom is defined externally.

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5 Generalized atoms in [2] are essentially conjunctions of generalized atoms defined here.

6 Note that domain plays no role in semantics, but is needed for complexity measures. In aggregates, multi-sets are allowed. In this case, occurrences of the same atom in a multi-set can be represented by different atoms whose equivalences can be enforced by additional rules.
2.2 Logic programs with generalized atoms

A logic program with generalized atoms (or sometimes just called a program) is a finite set of rules of the form: \( a \leftarrow \beta_1, ..., \beta_m, \neg \beta_{m+1}, ..., \neg \beta_n \), where \( m, n \geq 0 \), \( a \) is an ordinary atom and \( \beta_i \) (\( 1 \leq i \leq n \)) are generalized atoms. A rule is normal if all \( \beta_i \) are ordinary atoms, and a program is normal if all rules in it are normal. A program is called a logic program with monotone and anti-monotone generalized atoms, if every generalized atom in \( P \) is monotone or anti-monotone.

For a rule \( r \) of the above form, the head of the rule is denoted by \( H(r) = a \) and the body of the rule by \( B(r) = \{ \beta_1, ..., \beta_m, \neg \beta_{m+1}, ..., \neg \beta_n \} \). Also, we define \( \text{Pos}(r) = \{ \beta_1, ..., \beta_m \} \) and \( \text{Neg}(r) = \{ \neg \beta_{m+1}, ..., \neg \beta_n \} \) to denote positive atoms and negative atoms of \( B(r) \) respectively. We may use sets as conjunctions. A generalized literal is either a generalized atom \( \alpha \), or its negation, \( \neg \alpha \). Note that, without confusion, ordinary literals are (special cases of) generalized literals.

Let \( I \) be an interpretation, \( \beta \) a generalized atom and \( r \) a rule. Recall that if \( \beta \) maps \( I \) to \( t \) we say that \( I \) satisfies \( \beta \) and write \( I \models \beta \). We define that \( I \models \neg \beta \) if \( \beta \) maps \( I \) to \( f \), and \( I \models B(r) \) if \( I \models l \) for all \( l \in B(r) \). \( I \) is a model of a set of rules \( P \) if \( I \) satisfies every \( r \in P \).

Let \( I \) be a set of atoms. Then by \( \bar{I} = \{ a \mid a \in \Sigma \setminus I \} \), we denote the set of atoms exclusive of \( I \).

Well-founded semantics are typically defined by building a partial interpretation. Let \( S \) be a set of literals. We define \( S^+ = \{ a \mid a \in S \} \), \( S^- = \{ a \mid \neg a \in S \} \), and \( \neg S = \{ \neg a \mid a \in S \} \). \( S \) is consistent if \( S^+ \cap S^- = \emptyset \). A partial interpretation \( S \) is a consistent subset of \( \Sigma \cup \neg \Sigma \). Any atom not appearing in \( S \) is said to be undefined. A consistent extension of \( S \) is an interpretation \( I \) such that \( S \subseteq I \cup \neg I \). Note that in this paper we only consider consistent extensions that are total. In the sequel, we restrict \( \Sigma \) to the set of atoms appearing in \( P \). This is typically called the Herbrand base of \( P \) and is denoted by \( \text{HB}_P \).

2.3 Well-founded semantics for normal logic programs

To place our work in perspective, we briefly review the well-founded semantics for normal logic programs, which can be defined alternatively in several ways, one of which is based on the notion of unfounded set, which we adopt here.

Let \( P \) be a normal logic program and \( S \) a partial interpretation. A set \( U \subseteq \Sigma \) is an unfounded set of \( P \) w.r.t. \( S \), if for every \( a \in U \) and every rule \( r \in P \) with \( H(r) = a \), either (i) \( \neg b \in S \cup \neg U \) for some \( b \in \text{Pos}(r) \), or (ii) \( b \in S \) for some \( b \in \text{Neg}(r) \). The greatest unfounded set of \( P \) w.r.t. \( S \) exists, which is denoted by \( U_P(S) \).

Intuitively, unfounded atoms are those that can be safely assumed to be false without affecting the evaluation of the rules under the given interpretation.
We then define two operators

\[ T_P(S) = \{ H(r) \mid r \in P, \text{Pos}(r) \cup \neg \text{Neg}(r) \subseteq S \} \]

\[ W_P(S) = T_P(S) \cup \neg U_P(S) \]

The operator \( W_P \) is monotone, and thus has a least fixpoint, which defines the well-founded semantics of \( P \).

3 Well-Founded Semantics for Logic Programs with Generalized Atoms

In the definition of the well-founded semantics for normal logic programs, when an atom \( a \) is in a partial interpretation \( S \), it is clear that \( a \) remains to be satisfied by all consistent extensions of \( S \). However, this is not the case in general for a non-monotone generalized atom. We thus extend the notion of truth (resp. falsity) to the notion of persistent truth (resp. persistent falsity) under a partial interpretation.

Definition 1. Let \( \alpha \) be a generalized atom and \( S \) a partial interpretation.

- if \( \alpha \) is an ordinary atom, it is persistently true (resp. persistently false) under \( S \) if \( \alpha \in S^+ \) (resp. \( \alpha \in S^- \));
- Otherwise, \( \alpha \) is persistently true (resp. persistently false) under \( S \) if for all consistent extensions \( I \) of \( S \), \( I \models \alpha \) (resp. \( I \not\models \alpha \)).
- not \( \alpha \) is persistently true under \( S \) if \( \alpha \) is persistently false under \( S \).
- not \( \alpha \) is persistently false under \( S \) if \( \alpha \) is persistently true under \( S \).

Intuitively, a generalized atom \( \alpha \) is persistently true (resp. persistently false) relative to a partial interpretation \( S \), iff \( \alpha \) is true (resp. false) relative to all consistent extensions of \( S \), i.e., the truth or falsity of \( \alpha \) remains unaffected when all undefined atoms w.r.t. \( S \) get assigned in any possible way.

The definition above naturally extends to conjunctions of generalized literals. The following definition can be seen as a paraphrase of the notion of unfounded set for logic programs with aggregates [12].

Definition 2. (Unfounded set) Let \( P \) be a logic program with generalized atoms and \( S \) a partial interpretation. A set \( U \subseteq HB_P \) is an unfounded set of \( P \) relative to \( S \) if for each \( r \in P \) with \( H(r) \in U \), some generalized literal in \( B(r) \) is persistently false w.r.t. \((S \setminus U) \cup \neg U\).

The definition says that an atom \( a \) is in an unfounded set \( U \), relative to \( S \), because, for every rule with \( a \) in the head, at least one body literal is false in all consistent extensions of \((S \setminus U) \cup \neg U\).

Definition 3. (Unfounded-free interpretation) Let \( S \) be a partial interpretation of a logic program with generalized atoms \( P \). \( S \) is unfounded-free, if \( S \cap U = \emptyset \), for each unfounded set \( U \) of \( P \) relative to \( S \).

The following lemma has been proved for logic programs with aggregates [12]. That proof can be adopted easily to programs with generalized atoms.
Lemma 1. Let \( P \) be a logic program with generalized atoms and \( S \) be an unfounded-free interpretation. (i) Unfounded sets of \( P \) relative to \( S \) are closed under union. (ii) The greatest unfounded set of \( P \) relative to \( S \) exists, which is the union of all unfounded sets of \( P \) relative to \( S \).

Now we define the well-founded semantics. The following operators are needed for the definition of well-founded semantics.

Definition 4. Let \( P \) be a logic program with generalized atoms and \( S \) an unfounded free partial interpretation. Define the operators \( T_P, U_P, \) and \( W_P \) as follows:

(i) \( T_P(S) = \{ H(r) \mid r \in P, B(r) \text{ is persistently true under } S \} \).
(ii) \( U_P(S) \) is the greatest unfounded set of \( P \) relative to \( S \).
(iii) \( W_P(S) = T_P(S) \cup \neg U_P(S) \).

As a notation, we define \( W^0_P = \emptyset \), and \( W^{i+1}_P = W_P(W^i_P) \), for all \( i \geq 0 \). Note that \( W^0_P \) is an unfounded-free interpretation and so is every \( W^i_P \), for \( i \geq 0 \). Thus in every step the greatest unfounded set is computed relative to an unfounded-free set.

Lemma 2. The operators \( T_P, U_P, \) and \( W_P \) are all monotone.

Note that the operator \( W_P \) is well-defined, i.e., it is a mapping on unfounded-free partial interpretations, which, along with the subset relation, forms a complete lattice. Since \( W_P \) is monotone, it follows from the Knaster-Tarski fixpoint theorem that the least fixpoint of \( W_P \) exists [25].

Definition 5. (Well-founded semantics) Let \( P \) be a logic program with generalized atoms. The well-founded semantics (WFS) of \( P \), denoted \( \text{WFS}(P) \), is defined as the least fixpoint of the operator \( W_P \), denoted \( \text{lfp}(W_P) \). An atom \( a \in \Sigma \) is well-founded (resp. unfounded) relative to \( P \) iff \( a \) (resp. \( \neg a \)) is in \( \text{lfp}(W_P) \).

Observe that, the only difference in the operators defined here from those defined for normal logic programs is in the evaluation of body literals - being true (resp. false), has been replaced by being persistently true (resp. persistently false). It thus follows that the well-founded semantics for logic programs with generalized atoms is a generalization of the well-founded semantics for normal logic programs, and it treats monotone, anti-monotone, and non-monotone generalized atoms in a uniform manner.

Example 2. Consider program \( P \) below, where generalized atoms are aggregates.

\[
\begin{array}{ll}
  r_1 & : p(-1). \\
  r_2 & : p(-2) \leftarrow \text{SUM}(\{ X \mid p(X) \}) \leq 2. \\
  r_3 & : p(3) \leftarrow \text{SUM}(\{ X \mid p(X) \}) > -4. \\
  r_4 & : p(-4) \leftarrow \text{SUM}(\{ X \mid p(X) \}) \leq 0.
\end{array}
\]

The aggregates under \( \text{SUM} \) are self-explaining, e.g., \( \text{SUM}(\{ X \mid p(X) \}) \leq 2 \) means that the sum of \( X \) for satisfied atoms \( p(X) \) is less than or equal to 2. We start with \( W^0_P = \emptyset \), and then \( W^1_P = \{ p(-1) \} \). Observe that the body of \( r_2 \) is persistently true under \( W^1_P \).

We then have \( W^2_P = \{ p(-1), p(-2) \} \), and \( W^3_P = \{ p(-1), p(-2), p(-4) \} \). Now the body of \( r_4 \) is persistently false under \( W^3_P \cup \{ \neg p(3) \} \). So, \( W^4_P = \{ p(-1), p(-2), p(-4), \neg p(3) \} \) and, \( \text{WFS}(P) = \text{lfp}(W_P) = W^4_P \).
3.1 Complexity

Here, let us assume that for a generalized atom \( \alpha \), \( \text{Dom}(\alpha) \) is finite and the relation \( I \models \alpha \) (resp. \( I \models \neg \alpha \)) can be determined in polynomial time in the size of \( \text{Dom}(\alpha) \). This is the case for practical aggregates in logic programming, reasoning with Horn clauses, and satisfiability testing of the combined complexity for the DL-Lite family of DLs [4].

In general, the problem of computing the WFS of a program is intractable since determining whether a generalized atom is persistently true or persistently false under a partial interpretation is in general intractable.

**Proposition 1.** Let \( \alpha \) be a generalized atom. Checking whether \( \alpha \) is persistently true relevant to a partial interpretation \( S \) is Co-NP-complete.

However, the WFS for logic programs with monotone and anti-monotone generalized atoms is tractable. Let \( \alpha \) be a generalized atom and \( S \) a partial interpretation. To check whether \( \alpha \) is persistently true under \( S \), if \( \alpha \) is monotone we test \( S^+ \models \alpha \), and if \( \alpha \) is anti-monotone we test \( S^- \not\models \alpha \). On the other hand, if \( \alpha \) is monotone, then \( \alpha \) is persistently false under \( S \) iff \( S^- \not\models \alpha \), and if \( \alpha \) is anti-monotone, then \( \alpha \) is persistently false under \( S \) iff \( S^+ \not\models \alpha \).

As the number of distinct atoms is at most the size of \( P \), and the greatest unfounded set can be generated incrementally in polynomial time, the following proposition can be proved in a way similar to the claim for logic programs with monotone and anti-monotone aggregates.

**Proposition 2.** Let \( P \) be a logic program with monotone and anti-monotone generalized atoms. Then, \( \text{lfp}(W_P) \) can be computed in polynomial time.

4 Polynomial Approximation for DL-Programs

4.1 Description Logic Program

We assume that the reader has some familiarity with description logics (DLs) [5], which are decidable fragments of first order logic.

A dl-program is a combined knowledge base \( \mathcal{KB} = (L, P) \), where \( L \) is a DL knowledge base, which is a collection of axioms in the underlying DL, and \( P \) is a rule base, which is a finite set of rules of this form: \( h \leftarrow A_1, \ldots, A_m, \neg B_1, \ldots, \neg B_n \), where \( h \) is an atom, and \( A_i \) and \( B_i \) are atoms or dl-atoms,\(^7\) which are of the form

\[
DL[S_1 op_1 p_1, \ldots, S_m op_m p_m; Q](t)
\]

in which \( S_i \) is a concept or role from the vocabulary of \( L \), \( op_i \in \{\cup, \cap\} \), \( p_i \) is a predicate symbol only appearing in \( P \) whose arity matches that of \( S_i \), and \( Q(t) \) is called a dl-query,\(^8\) where \( t \) is a list of constants (in the following, we use bold letters

---

\(^7\) For simplicity, we assume that equality doesn’t appear in rules.

\(^8\) A dl-query is of the form \( Q(t) \), where \( t \) is a list of terms, and \( Q \) is a concept, a role, or a concept inclusion axiom, built from the vocabulary of \( L \).
for this purpose), and \( Q \) is a concept, a role, or a concept inclusion axiom, built from the vocabulary of \( L \).

Intuitively, \( S_i \cup p_i \) extends \( S_i \) by the extension of \( p_i \), and \( S_i \cup p_i \) analogously extends \( \neg S_i \); the expression \( S \cap p \) instead constrains \( S_i \) to \( p_i \). It is clear that the operator \( \cap \) (which we call the constraint operator) may cause a dl-program to be non-monotone; a dl-atom which is free of the constraint operator is monotone.

In this paper, we assume that a dl-program is ground, and define the Herbrand base of \( P \), denoted \( HB_P \), to be the set of all ground atoms \( p(t) \), where \( p \) appears in \( P \) and \( t \) is a tuple of constants. Interpretations are subsets of \( HB_P \).

**Definition 6.** Let \( KB = (L, P) \) be a dl-program and \( I \subseteq HB_P \) an interpretation. We define the satisfaction relation under \( L \), denoted \( \models_L \), as follows:

1. \( I \models_L \top \) and \( I \not\models_L \bot \).
2. For any atom \( a \in HB_P \), \( I \models_L a \) if \( a \in I \).
3. For any (ground) dl-atom \( A = DL[S_1 op_1 p_1, \ldots, S_m op_m p_m; Q](e) \) occurring in \( P \), \( I \models_L A \) if \( L \cup \bigcup_{i=1}^{m} A_i \models Q(e) \), where

\[
A_i = \begin{cases} 
\{ S_i(e) \mid p_i(e) \in I \}, & \text{if } op_i = \cup; \\
\{ \neg S_i(e) \mid p_i(e) \in I \}, & \text{if } op_i = \cap; \\
\{ \neg S_i(e) \mid p_i(e) \notin I \}, & \text{if } op_i = \lor.
\end{cases}
\]

We now define a mapping \( \xi \) from dl-programs to programs with generalized atoms.

**Definition 7.** Let \( KB = (L, P) \) be a dl-program, and \( \xi(KB) \) denote the set of rules obtained from \( KB \) in the following way: for each \( r \in P \), we have \( r' \in \xi(KB) \) such that \( r' \) is \( r \) except that each dl-atom \( q \) appearing in \( r \) is replaced by a generalized atom \( \alpha_q \) such that

- \( \text{Dom}(\alpha_q) \) is the set of atoms appearing in \( P \), and
- We identify \( \Sigma \) with \( HB_P \), and define the semantics of \( \alpha_q \) as: for all interpretations \( I, I \models_L q \) iff \( I \models \alpha_q \).

That is, \( \xi \) is a transformation that preserves satisfiability under \( \models_L \).

**Definition 8.** Let \( KB = (L, P) \) be a dl-program. The well-founded semantics (WFS) of \( KB \) is defined in terms of the well-founded semantics of \( \xi(KB) \).

**Example 3.** Consider a dl-program \( KB = (\emptyset, P) \), where \( P \) consists of

\[
\begin{align*}
  r_1 : & \quad p(a) \not\rightarrow \text{not } DL[S_1 \cap q, S_2 \cup r; S_1 \cap S_2](a). \\
  r_2 : & \quad q(a) \not\rightarrow DL[S \cap q; S](a). \\
  r_3 : & \quad r(a) \not\rightarrow DL[S \cap q; \neg S](a).
\end{align*}
\]

Let \( P' = \xi(KB) \). We have \( W^0_P = \emptyset \) and, for example, we do not derive \( p(a) \) since there is a consistent extension of \( \emptyset \) that satisfies the dl-atom in rule \( r_1 \), but \( \{ q(a) \} \) is an unfounded set relative to \( \emptyset \). The reader can verify that \( W^1_P = \{ \neg q(a) \} \), \( W^2_P = \{ \neg q(a), r(a) \} \), and \( W^3_P = \{ \neg q(a), \neg p(a), r(a) \} \), which is the least fixpoint of \( W_P \).

Due to the one-to-one correspondence between satisfaction under \( \models_L \) and satisfaction under \( \models \), in the following, we may refer to the WFS of a dl-program directly, without going through the mapping \( \xi \). Thus, the WFS of a dl-program \( KB \) is referred to as \( \text{WFS}(KB) \). We therefore can apply the notion of persistent truth and falsity directly to dl-programs.
4.2 Removing non-monotone dl-atoms as a polynomial approximation

Clearly, non-monotone dl-atoms are the result of applying the constraint operator \( \land \), as the satisfiability of such atoms depends on truth value of propositional atoms exclusive of a given interpretation. In [11], the authors suggest a polynomial time rewrite to remove the constraint operator. The idea is to encode a dl-atom that contains the operator \( \land \) in terms of the operator \( \lor \), and thus the resulting dl-program is free of \( \land \). It is known that these dl-atoms are monotone; then for such a dl-atom \( \alpha \), \( \neg \alpha \) is anti-monotone. Hence the resulting dl-program can be viewed as a program with monotone and anti-monotone generalized atoms, and its WFS can be computed in polynomial time.

Definition 9. (Transformation \( \pi \)) Let \( KB = (L, P) \) be a dl-program. \( \pi(KB) \) is a dl-program obtained from \( KB \) by

1. replacing each occurrence of \( S_i \land p_i \) with \( S_i \lor \bar{p}_i \) (\( \bar{p}_i \) is the complement of \( p_i \)), and
2. for each atom \( p_i(t) \) that occurs in the head of some rule in \( P \), add the following rule

\[
\bar{p}_i(t) \leftarrow \text{notDL}[S'_i \lor p_i; S'_i](t)
\]

where \( S'_i \) is a fresh concept or role name.

Note that \( \pi \) is a polynomial transformation.

Since the transformation \( \pi \) does not affect DL knowledge base \( L \) in \( KB \), we may write \( \pi(P) \) to denote the set of all transformed rules, and \( \pi(r) \) to denote the set of transformed rules for \( r \in P \). By the transformation \( \pi \), any dl-program \( KB = (L, P) \) can be rewritten to a dl-program \( \pi(KB) \) free of the constraint operator, and thus \( \pi(KB) \) only contains monotone and anti-monotone dl-atoms. By Proposition 2, we know that, if query answering with the underlying DL is tractable\(^9\), then the WFS of \( \pi(KB) \) can be computed in polynomial time.

The question is whether the transformation \( \pi \) is faithful. That is, whether it is the case that for all dl-programs \( KB \), the WFS of \( KB \) is equivalent to the WFS of \( \pi(KB) \), barring freshly added concept/role names in \( \pi(KB) \). The following example shows that the answer to this question is no.

Example 4. Consider a single-rule dl-program \( KB = (\emptyset, P) \), where \( P \) consists of a single rule

\[
p(a) \leftarrow \text{DL}[S \lor p, S \land p; \neg S](a).
\]

The WFS of \( KB \) is \( \{p(a)\} \). On the other hand, \( \pi(P) \) consists of two rules

\[
p(a) \leftarrow \text{DL}[S \lor p, S \lor \bar{p}; \neg S](a).
\bar{p}(a) \leftarrow \text{notDL}[S'_i \lor p_i; S'_i](a).
\]

Clearly, the WFS of \( \pi(KB) \) is \( \emptyset \), as neither \( p(a) \) nor \( \bar{p}(a) \) is derivable under \( \emptyset \), and neither is unfounded relative to \( \emptyset \).

\(^9\) Note that this requires that not only the underlying DL is tractable, but also query in a dl-atom is bounded.
However, we show that the transformation $\pi$ is correct, i.e., given a dl-program $KB$, when restricted to the language of $KB$, all the well-founded (resp. unfounded) atoms relative to $\pi(KB)$ are well-founded (resp. unfounded) relative to $KB$.

Let $KB = (L, P)$ be a dl-program and $\tau(KB)$ be its signature. We denote the WFS of $\pi(KB)$, restricted to $\tau(KB)$, by $WFS(\pi(KB))|_{\tau(KB)}$ (similarly, $lfp(WFS(\pi(KB)))|_{\tau(KB)}$).

For simplicity, we just write $\tau$ for $\tau(KB)$.

We now give the main theorem of this section.

**Theorem 1.** Let $KB = (L, P)$ be a dl-program. Then $WFS(\pi(KB))|_{\tau} \subseteq WFS(KB)$.

Please see the appendix for the proof of Theorem 1.

## 5 Polynomial Approximation for Logic Programs with Aggregates

When generalized atoms are aggregates, logic programs with generalized atoms are logic programs with aggregates. The mapping from logic program with aggregates to logic program with generalized atoms is similar to the mapping $\xi$, which can be obtained by replacing dl-atom in Definition 7 with aggregate atom. In this section, we present a case study on polynomial approximation for logic programs with standard aggregates [23].

### 5.1 Syntax of logic programs aggregates

Following [23], an aggregate (or aggregate atom) is a constraint on a set of atoms taking the form

$$aggr\{\{X \mid p(X)\}\} op Result$$

where $aggr$ is an aggregate function. The standard aggregate functions are those in $\{SUM, COUNT, AVG, MAX, MIN\}$. The set $\{X\mid p(X)\}$ is called an intensional set, where $p$ is a predicate, and $X$ is a variable which takes value from a set $D(X) = \{a_1, \ldots, a_n\}$, called the variable domain. The relational operator $op$ is from $\{=, \neq, <, >, \leq, \geq\}$ and $Result$ is a numeric constant.

The domain of an aggregate $A$, denoted $Dom(A)$, is the set of atoms $\{p(a) \mid a \in D(X)\}$. The size of an aggregate is $|Dom(A)|$.

For an aggregate $A$, the intensional set $\{X\mid p(X)\}$, the variable domain $D(X)$, and the domain of an aggregate $Dom(A)$ can also be a multiset which may contain duplicate members.

Let $M$ be a set or multiset of atoms. $M$ is a model of (satisfies) an aggregate $A$, denoted $M \models A$, if $aggr\{\{a \mid p(a) \in M \cap Dom(A)\}\} op Result$ holds, otherwise $M$ is not a model of (does not satisfy) $A$, denoted $M \not\models A$.

For instance, consider the aggregate $A = SUM(\{X\mid p(X)\}) \geq 2$, where $D(X) = \{-1, 1, 2\}$. For the sets $M_1 = \{p(2)\}$ and $M_2 = \{p(-1), p(1)\}$, we have $M_1 \models A$ and $M_2 \not\models A$. For the multiset $M_3 = \{p(1), p(1)\}$, we have $M_3 \models A$.

A logic program with aggregates (or aggregate logic program) is a finite set of rules of this form: $h \leftarrow A_1, \ldots, A_k, notB_1, \ldots, notB_{m}, G_1, \ldots, G_n$, where $h, A_i$ and $B_j$ are atoms and $G_i$ are aggregates.
5.2 Disjunctive rewriting as a polynomial approximation

To optimize programs with constraint atoms, in [15], replacement techniques are studied, where a complex constraint may be decomposed into simpler ones. In one replacement scheme, the authors propose to rewrite a program with disjunctive encoding for c-atoms under the answer set semantics. The idea is to encode a complex c-atom by a disjunction of simpler c-atoms. We apply this idea to aggregates.

A disjunctive encoding of an aggregate $A$ is a disjunction of aggregates $A_i$ ($1 \leq i \leq m$), denoted by $d(A_1, \ldots, A_m)$, such that for any interpretation $I$, $I \models A$ iff $I \models d(A_1, \ldots, A_m)$. That is, disjunctive encoding preserves satisfaction.

In [23], the authors show that the determination of persistent truth of an aggregate atom involving $\text{SUM}/\text{AVG}$ and $\neq$ at the same time is intractable, while determining the same for all other aggregate atoms is tractable. Now, by definition, an aggregate atom $A$ is persistently false under $S$ iff the complement of $A$ is persistently true under $S$. As a result, determining persistent falsity of $\text{SUM}/\text{AVG}$ involving the $=$ operator is also intractable. Theorem 2 gives a sufficient condition for computing the WFS of an aggregate logic program in polynomial time.

Theorem 2. Let $P$ be an aggregate logic program. In the computation of its WFS, if there is no need to compute persistent truth for any aggregate atom that involves $\text{SUM}/\text{AVG}$ and the operator $\neq$, and there is no need to compute persistent falsity of any aggregate atom that involves $\text{SUM}/\text{AVG}$ and the operator $=$, then the computation of $WFS(P)$ is a polynomial process.

Thus the goal of disjunctive rewriting for aggregate logic programs is to transform away aggregates of the form $f(\cdot) \neq c$ for computing well-founded atoms and $f(\cdot) = c$ for computing unfounded atoms, where $f \in \{\text{SUM, AVG}\}$.

Definition 10. (Disjunctive rewriting) Let $P$ be a logic program with aggregate. The disjunctive rewrite of $P$ produces two programs, one for polynomial time computation of well-founded atoms of $P$, denoted as $P_w$ and the other for polynomial time computation of unfounded atoms of $P$, denoted as $P_u$. We define $P_w$ as:

For each occurrence of aggregate atom of the form $f(\cdot) \neq c$ in $P$, where $f \in \{\text{SUM, AVG}\}$, we replace that atom with a unique symbol $\alpha$ and add the following two rules: $\alpha \leftarrow f(\cdot) > c$ and $\alpha \leftarrow f(\cdot) < c$.

and define $P_u$ as:

For each occurrence of aggregate atom of the form $f(\cdot) = c$ in $P$, where $f \in \{\text{SUM, AVG}\}$, we replace that atom with the conjunction of two aggregates, $f(\cdot) \leq c$ and $f(\cdot) \geq c$.

By an abuse of notation, let us denote the pair of programs $P_w$ and $P_u$ by $P_{(w,u)}$. Now, we revise the definition of the operator $W_P$ (cf. Definition 4) as follows:

$$W_{P_{(w,u)}}(S) = TP_u(S) \cup \neg UP_w(S).$$

It can be shown that the operator $W_{P_{(w,u)}}$ is monotone, thus its least fixpoint can be computed iteratively. Again, let us denote by $WFS(P_{(w,u)})$ the least fixpoint of the operator $W_{P_{(w,u)}}$.

The following example shows that the disjunctive rewrite is an incomplete method.
Example 5. Consider the following aggregate logic program $P$.

$$
\begin{align*}
  p(2) & \leftarrow \text{SUM}([X \mid p(X)]) \neq -1. \\
p(-3) & \leftarrow p(2). \\
  p(1) & \leftarrow .
\end{align*}
$$

WFS($P$) is computed as follows: $W^P_S = \emptyset, W^P_p = \{p(1)\}, W^P = \{p(1), p(2)\}$, and $W^P_{\emptyset} = \{p(1), p(2), p(-3)\}$, which is the least fixpoint of $W_P$.

By disjunctive rewriting, we have $P_w$ below

$$
\begin{align*}
p(2) & \leftarrow \alpha, \\
p(-3) & \leftarrow p(2). \\
\alpha & \leftarrow \text{SUM}([X \mid p(X)]) > -1. \\
p(1) & \leftarrow .
\end{align*}
$$

As $P$ does not contain any aggregate atom of the form $f(.) = c, P_w = P$. Now starting with the interpretation $\emptyset$, $T_{P_w}(\emptyset) = \{p(1)\}$ and $U_{P_w}(\emptyset) = \emptyset$, thus $W_{P_{w,u}}(\emptyset) = \{p(1)\}$. It can be verified easily that this is a fixpoint of $W_{P_{w,u}}, i.e., WFS(P_{w,u}) = \{p(1)\}$.

It can be seen that disjunctive rewriting produces stronger constraints. Let us extend the notion of persistent truth and falsity to disjunction of literals in a natural way. Then, we can see that, given a disjunction of aggregates, say $(f(.) < c) \lor (f(.) > c)$, and a partial interpretation $S$, the fact that $f(.) < c$ is persistently true under $S$ or $f(.) > c$ is persistently true under $S$ implies that $(f(.) < c) \lor (f(.) > c)$ is persistently true under $S$, but the converse does not hold in general - that $(f(.) < c) \lor (f(.) > c)$ is satisfied by all consistent extensions of $S$ does not imply that $f(.) < c$ is persistently true under $S$ or $f(.) > c$ is persistently true under $S$, because it may be due to that some consistent extensions satisfy $f(.) < c$ and others satisfy $f(.) > c$. Similar for the computation of unfounded set, that $f(.) = c$ is persistently false iff $f(.) \neq c$ is persistently true if either $f(.) > c$ is persistently true or $f(.) < c$ is persistently true iff either $f(.) \leq c$ is persistently false or $f(.) \geq c$ is persistently false. The converse for the if statement above does not hold because that $f(.) = c$ is persistently false may be due to the fact that for some consistent extensions $f(.) < c$ holds and for the others $f(.) > c$ holds.

The above arguments actually give a proof sketch of the following theorem.

Theorem 3. Let $P$ be a logic program with aggregates. Then, $\text{WFS}(P_{w,u}) \subseteq \text{WFS}(P)$.

6 Related Work and Discussion

The well-founded semantics defined in this paper for logic programs with generalized atoms is based on essentially the same notion of unfounded set formulated in [12]. By the work of [20], for logic programs with aggregates, this well-founded semantics is known to approximate answer sets based on the notion of conditional satisfaction [24]. This well-founded semantics is different from that of [27], which approximates answer sets by reduct [24]. The WFS of [27] is weaker than the WFS defined here, but without any reduction on complexity.

In [14], the authors present a well-founded semantics for hybrid Minimal Knowledge and Negation as Failure (MKNF) knowledge bases, which is a local closed world.
extension of the MKNF DL knowledge base. The well-founded semantics defined in [14] is shown to be tractable, if the chosen DL fragment is tractable. As shown in Proposition 1 of this paper, even if we assume the entailment relation $I \models \varphi$ is tractable, for interpretation $I$ and generalized atom $\varphi$, computing the WFS is still not. This is inevitable since classic formulas under the scope of negation are anti-monotone while generalized atoms may be neither monotone nor anti-monotone. The precise relationship between the MKNF WFS and the WFS defined here requires further study.

If we assume that the domain of a generalized atom is bounded [3], the well-founded semantics can be computed in polynomial time. This assumption is reasonable only for generalized atoms with small domains.

For improving propagation efficiency for HEX-programs, a decision criterion is introduced in [8] to allow to decide if further check is necessary (with the external sources) to complete the computation of the Unfounded Set (UFS) of the guessing program $Q$ obtained from a given HEX-program $P$ w.r.t. an interpretation $A$. The decision criterion is as follows: is there any atom dependency cycle that exists in $P$, which contains external edges (e-cycle)? Following this decision criterion, the authors devise a program decomposition technique which decomposes a given HEX-program into two types of components - one type of component is with e-cycles and other type of component does not have e-cycles. UFS checking is needed only for the components which do have e-cycles. Thus this technique avoids UFS checking when it is not necessary. This work however does not prevent complexity jump in constraint propagation. The decision criterion reduces computational cost linearly. In our case studies, we focus on subtle aspects of computation that cause complexity jump, which may be avoided by incomplete methods. We establish the links between such incomplete methods with the well-founded semantics of the underlying logic program.

We wonder whether the idea of incomplete methods can be pursued for HEX-programs in general. If yes, it will be interesting to study the characterizations of the type of information that may be lost.

References


A Appendix : Proof of Theorem 1

Proof. We prove the claim by induction on the construction of lfp($W_{\pi(KB)}^k$) and lfp($W_{KB}$).

(a) Base: $W_{\pi(KB)}^0|_\tau = W_{KB}^0|_\tau = \emptyset$.

(b) Step: Assume, for all $k \geq 0$, $W_{\pi(KB)}^k|_\tau \subseteq W_{KB}^k$, and prove $W_{\pi(KB)}^{k+1}|_\tau \subseteq W_{KB}^{k+1}$.

By definition, we know

$$W_{\pi(KB)}^{k+1}|_\tau = T_{\pi(KB)}(W_{\pi(KB)}^k)|_\tau \cup \neg U_{\pi(KB)}(W_{\pi(KB)}^k)|_\tau$$

$$W_{KB}^{k+1} = T_{KB}(W_{KB}^k) \cup \neg U_{KB}(W_{KB}^k)$$

(3)

(4)

To prove $W_{\pi(KB)}^{k+1}|_\tau \subseteq W_{KB}^{k+1}$, it is sufficient to prove both of

$$T_{\pi(KB)}(W_{\pi(KB)}^k)|_\tau \subseteq T_{KB}(W_{KB}^k)$$

$$U_{\pi(KB)}(W_{\pi(KB)}^k)|_\tau \subseteq U_{KB}(W_{KB}^k)$$

(5)

(6)

Below, let us assume that at most one dl-atom may appear in a rule. The proof can be generalized to arbitrary rules by the same argument, for the transformation of one dl-atom at a time.

(i) We first prove (5). Let $a \in T_{\pi(KB)}(W_{\pi(KB)}^k)|_\tau$. By definition, $\exists r' \in \pi(P)$ such that $B(r')$ is persistently true under $W_{\pi(KB)}^k$. WLOG, assume for some $r \in P$, $\pi(r) = \{r', r''\}$, as illustrated in (7) below, in which a dl-atom appears positively, which is replaced by rules in (8) of $\pi(r)$.

$$r : a \leftarrow ..., DL[...], S_j \subseteq p_j, ..., Q](c), ...$$

$$r' : a \leftarrow ..., DL[...], S_j \cup p_j, ..., Q](c), ...$$

$$r'' : p_j(e) \leftarrow not DL[S_j \cup p_j; S'](e)$$

(7)

(8)

Let $D$ denote the dl-atom in (7), and $D'$ the corresponding dl-atom in (8). If the operator $\cap$ does not occur in $D$, then trivially $B(r')$ is persistently true under $W_{KB}^k$, and thus $a \in T_{KB}(W_{KB}^k)$.

Otherwise, since $B(r')$ is persistently true under $W_{\pi(KB)}^k$, we have $D'$ is persistently true under $W_{\pi(KB)}^k$, and by the fact that $D'$ is monotone, we have $W_{\pi(KB)}^k|_{D'} \equiv L D'$. The atom $p_j(e)$ may or may not play a role in the entailment $L \cup \bigcup_{i=1}^n D_i' \models Q(e)$ (cf. Definition (6)). The proof is trivial if it does not. Otherwise, $p_j(e)$ is well-founded already w.r.t. $W_{\pi(KB)}^k$, and by the last rule in (8), $p_j(e)$ is unfounded w.r.t. $W_{\pi(KB)}^k$, for some $k' < k$. By induction hypothesis, we know $W_{\pi(KB)}^k|_{D'} \subseteq W_{KB}^{k'}$, thus $p_j(e)$ is unfounded w.r.t. $W_{KB}^{k'}$. It follows $D$ is persistently true under $W_{KB}^k$, and by the assumption that $D$ is the only dl-atom in (7) and that $B(r')$ is persistently true under $W_{\pi(KB)}^k$, we have $B(r)$ is persistently true under $W_{KB}^k$. Thus, $a \in T_{KB}(W_{KB}^k)$.

If $D$ appears negatively in rule body, the proof is similar because, given a partial interpretation $S$, not $D$ is persistently true under $S$ iff $D$ is persistently false under $S$, and we just need to swap well-founded and unfounded in the argument above.

(ii) To prove (6), assume that $a \in U_{\pi(KB)}(W_{\pi(KB)}^k)|_\tau$ and we show $a \in U_{KB}(W_{KB}^k)$.

Consider the case of (7). WLOG, assume that $r$ is the only rule in $P$ with $a$ in the head.
If the fact $a \in U_{\pi(\text{KB})}(W^k_{\pi(\text{KB})})$ is independent of $\bar{p}_j(e)$, the proof is trivial. Otherwise, that $a \in U_{\pi(\text{KB})}(W^k_{\pi(\text{KB})})$ is because $D'$ is persistently false under $\neg U' \cup W^k_{\pi(\text{KB})}$, where $U'$ is the greatest unfounded set relative to $W^k_{\pi(\text{KB})}$. This implies that $\bar{p}_j(e)$ must be unfounded w.r.t. $W^k_{\pi(\text{KB})}$, and it follows that $p_j(e)$ is well-founded already w.r.t. $W^{k'}_{\pi(\text{KB})}$, for some $k' < k$. Then, by induction hypothesis, we have that $p_j(e)$ is well-founded w.r.t. $W^k_{\text{KB}}$, which implies that $B(r)$ is persistently false under $\neg U \cup W^k_{\text{KB}}$, where $U$ is the greatest unfounded set w.r.t. $W^k_{\text{KB}}$. It follows $a \in U_{\text{KB}}(W^k_{\text{KB}})$. The proof for the case where a dl-atom appears negatively in rule body is similar.

Hence, the proof is completed. $\Box$
A Formal Semantics for the Cognitive Architecture ACT-R

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Abstract. The cognitive architecture ACT-R is very popular in cognitive sciences. It merges well-investigated results of psychology to a unified model of cognition. This enables researchers to implement and execute domain-specific cognitive models. ACT-R is implemented as a production rule system. Although its underlying psychological theory has been investigated in many psychological experiments, ACT-R lacks a formal definition from a mathematical-computational point of view.

In this paper, we present a formalization of ACT-R’s fundamental concepts including an operational semantics of the core features of its production rule system. The semantics abstracts from technical artifacts of the implementation. Due to its abstract formulation, the semantics is eligible for analysis. To the best of our knowledge, this operational semantics is the first of its kind.

Furthermore, we show a formal translation of ACT-R production rules to Constraint Handling Rules (CHR) and prove soundness and completeness of the translation mechanism according to our operational semantics.

Keywords: computational psychology, cognitive systems, ACT-R, production rule systems, Constraint Handling Rules, operational semantics

1 Introduction

Computational psychology is a field at the interface of psychology and computer science. It explores human cognition by implementing detailed computational models. The models are executable and hence capable of simulating human behavior. This enables researchers to conduct the same experiments with humans and a computational model to verify the behavior of the model. This leads to a gradual improvement of cognitive models. Furthermore, due to their executability, computational models have to be defined precisely. Hence, ambiguities which often appear in verbal-conceptual models can be eliminated.

Cognitive Architectures support the modeling process by bundling well-investigated research results from several disciplines of psychology to a unified theory. Domain-specific models are built upon such cognitive architectures. Ideally, cognitive architectures constrain modeling to only plausible domain-specific models.
Adaptive Control of Thought – Rational (ACT-R) is one of the most popular cognitive architectures \cite{2}. It is implemented as a production rule system. Although its underlying psychological theory is well-investigated and verified in many psychological experiments, ACT-R lacks a formal definition of its production rule system from a mathematical-computational point of view. I.e. the main data structures and the resulting operational semantics suggested by the psychological theory are not defined properly. This led to a reference implementation full of assumptions and technical artifacts beyond the theory making it difficult to overlook. Furthermore, the lack of a formal operational semantics inhibits analysis of the models like termination or confluence analysis.

In this paper, we present a formalization of the fundamental concepts of ACT-R leading to an operational semantics. The semantics abstracts from many details and artifacts of the implementation. Additionally, aspects like time or external modules are ignored to concentrate on the basic state transitions of the ACT-R production rule system. Those abstractions lead to a short and concise definition of the semantics making it suitable for theoretical analysis of the main aspects of ACT-R. Nevertheless, the semantics is still closely related to the general computation process of ACT-R implementations as we exemplify by a simple model.

The formalization of ACT-R has been derived from an implementation using Constraint Handling Rules (CHR), but captures the essence of the ACT-R production rule system. In this paper, we define the translation from ACT-R production rules to CHR rules formally. The translation is closely related to the translation process described informally in \cite{6}. Finally, we prove soundness and completeness of our translation according to our abstract operational semantics.

The paper is structured as follows: Section 2 covers the preliminaries. In section 3, we first recapitulate the formalization of the basic notions of ACT-R and then present its abstract operational semantics. The formal translation of ACT-R to Constraint Handling Rules is shown in section 4. Then, the translation is proved to be sound and complete in relation to our abstract operational semantics in section 5. In section 6, we present an example to show how our semantics works on a practical model. We conclude in section 7.

2 Preliminaries

First, we want to cover some notational aspects and introduce ACT-R and CHR.

2.1 Notation

We assume some basic notions of first-order logic and logic programming like syntactic equality or unification. Unification is denoted by $t[x/y]$ where all occurrences of the variable $x$ in the term $t$ are replaced by the variable $y$. For the sake of brevity, we treat logical conjunctions as multi-sets and vice-versa at some points. I.e. we use (multi-)set operators on logical conjunctions or multi-sets of terms as conjunctions. The symbol $\oplus$ is used to denote sequences of symbols. We use the relational notation for some functions and for binary relations we use an infix notation.
2.2 ACT-R

First of all, we describe ACT-R informally. For a detailed introduction to ACT-R we refer to [2] and [8]. Then we introduce a subset of its syntax and describe its operational semantics informally. The formalization is presented in section 3.

ACT-R is a production rule system which distinguishes two types of knowledge: 

- **declarative knowledge** holding static facts and 
- **procedural knowledge** representing processes which control human cognition.

For example, in a model of the game rock, paper, scissors, a declarative fact could be “The opponent played scissors”, whereas a procedural information could be that a round is won, if we played rock and the opponent played scissors. Declarative knowledge is represented as chunks. Each chunk consists of a symbolic name and labeled slots which hold symbolic values. The values can refer to other chunk names, i.e. chunks can be connected. Chunks are typed, i.e. the number and names of the slots provided by a chunk are determined by a type. As usual for production rule systems, procedural knowledge is represented as rules of the form IF conditions THEN actions. Conditions match values of chunks, actions modify them.

ACT-R has a modular architecture. For instance, there is a declarative module holding the declarative knowledge or a visual module perceiving the visual field and controlling visual attention. Each module has a set of affiliated buffers which can hold at most one chunk at a time. For example, there is a retrieval buffer which is associated to the declarative module and which holds the last retrieved chunk from the declarative memory.

The procedural system consists of a procedural memory with a set of production rules. The conditions of a production rule refer to the contents of the buffers, i.e. they match the values of the chunk’s slots.

There are three types of actions whose arguments are encoded as chunks as well: First of all, buffer modifications change the content of a buffer, i.e. the values of some of the slots of a chunk in a buffer. Secondly, the procedural module can state requests to external modules which then change the contents of buffers associated with them. Eventually, buffer clearings remove the chunk from a buffer. For the sake of simplicity, we only regard buffer modifications in this work. Additionally, we only refer to the procedural system with an arbitrary set of buffers which are not affiliated to a module. Nevertheless, our formalization and translation can be easily extended by other actions and modules [6]. Briefly, a request can be regarded as a modification, which modifies all information about a chunk according to a built-in function.

**Syntax** We define the syntax of an ACT-R production rule as follows:

**Definition 1 (production rule).** An ACT-R production rule \( r \) has the form

\[
(p \ x \ LHS \Rightarrow \ RHS)
\]

with

\[
\begin{align*}
LHS &= \bigoplus_{l=1}^{n_L} (b_l^L \ > \ \text{isa} \ t_l \ \bigoplus_{j=1}^{n_l} (s_{l,j}^L \ v_{l,j})) \\
RHS &= \bigoplus_{k=1}^{n_R} (b_k^R \ > \ \bigoplus_{i=1}^{n_i} (s_{k,i}^R \ v_{k,i}))
\end{align*}
\]
The \( b \) values are called buffers, the \( t \) values types, the \( s \) values slots and the \( v \) values are called values. The domain of the \( b, t \) and \( s \) values is \( S \), the set of symbols. The \( v \) values can be variables (denoted by the prefix \( = \)) or symbols. Variables and buffers appearing in RHS must also appear in LHS. However, slots on RHS are not required to appear on LHS. The buffers and slots of the RHS are assumed to be pairwise distinct. The slots on the RHS must be valid slots of the chunk type of the chunk in the buffer as specified on LHS \( \text{[3]} \).

Since the reference implementation of ACT-R is implemented in Lisp, the syntax of production rules reminds of the syntax of this language. Hence, a rule is simply a list of symbols (represented by the \( \oplus \) operator in our definition. A pair \( s_{i,j} v_{i,j} \) is called slot-value pair.

**Informal Operational Semantics** A production rule as defined in definition [1] is read as follows: The LHS of the rule are conditions matching the contents of the buffers. I.e. for a condition \( =b> s_1 v_1 s_2 v_2 \) the buffer \( b \) is checked for a chunk with the value \( v_1 \) in its \( s_1 \) slot and the value \( v_2 \) in its \( s_2 \) slot. If all conditions on the LHS match, the rule can be applied, i.e. the chunks in the buffers are modified according to the specification on the RHS. For an action \( =b> s_1 v'_1 \) the value in the slot \( s_1 \) of the chunk in buffer \( b \) is overwritten by the value \( v'_1 \). Since the buffers and slots on the RHS are pairwise distinct, there are no conflicting modifications.

**Example 1 (production rule).** Consider the following rule:

\[
(p \text{ recognize-win} \\
= b> \text{isa game me rock opponent scissors} \Rightarrow = b> \text{result win})
\]

It recognizes a win situation in the game rock, paper, scissors if the model has realized that the opponent played scissors and the agent played rock (which could be accomplished by a corresponding production rule interacting with the visual module). The situation is represented by a chunk of type game providing the slots me, opponent and result. As a result, it adds the information that the round has been won by modifying the result-slot of the buffer \( b \).

### 2.3 Constraint Handling Rules

We recap the syntax and semantics of Constraint Handling Rules (CHR) shortly. For a detailed introduction to the language, we refer to [4].

**Syntax** We briefly introduce a subset of the abstract CHR syntax as defined in [4]. Constraints are first-order logic predicates of the form \( c(t_1, \ldots, t_n) \) where the \( t \) values are terms. There are two distinct types of constraints: built-in and user-defined constraints. We constrain the allowed built-in constraints to true, false and the syntactic equality \( = \).
Definition 2 (CHR syntax). A CHR program $P$ is a finite set of rules. Simplification rules have the form

$$r @ H_k \setminus H_r \Leftrightarrow G \mid B.$$  

$r$ is an optional name of the rule, $H_k$ and $H_r$ are conjunctions of user-defined constraints (at least one of them is non-empty) called head constraints. $G$ is a conjunction of built-in constraints and is called the guard. Eventually, $B$ a conjunction of built-in and user-defined constraints and called the body of the rule.

Operational Semantics The operational semantics of CHR is defined as a state transition system. Hence, we first define the notion of a CHR state and then introduce the very abstract operational semantics of CHR [4] [5].

Definition 3 (CHR state). A CHR state is a goal, i.e. either true, false, a built-in constraint, a user-defined constraint or a conjunction of goals.

Definition 4 (head normal form). A CHR rule is in head normal form (HNF) if each argument of a head constraint is a unique variable.

A CHR rule can be put into HNF by replacing its head arguments $t_i$ with a new variable $V_i$ and adding the equations $V_i = t_i$ to its guard.

The operational semantics of CHR is defined upon a constraint theory $CT$ which is nonempty, consistent and complete and contains at least an axiomatization of the syntactic equality $=$ together with the built-in constraints $true$ and $false$.

Definition 5 (CHR operational semantics).

$$(H_k \land H_r \land G) \mapsto_r (H_k \land C \land H_k = H'_k \land H_r = H'_r \land B \land G)$$

if there is an instance with new variables $\bar{x}$ of a rule $r$ in HNF,

$$r @ H'_k \setminus H'_r \Leftrightarrow C \mid B.$$  

and $CT \models (G \rightarrow \exists \bar{x} \ (C \land (H_k = H'_k) \land (H_r = H'_r))).$

I.e., there is a state transition using the rule $r$, if there are constraints in the state which match the head constraints of the rule $r$ and the guard holds. For the successor state, the constraints in $H_k$ are kept, the constraints in $H_r$ are removed and the guard, body and the equality constraints from the matching are added.

3 Formalization of the ACT-R Production System

In this section, we formalize the core data structures of ACT-R formally. We follow the definitions from [6].
3.1 Chunk Stores

Intuitively, a chunk store represents a network of chunks. I.e., it contains a set of typed chunks. Each chunk has a fixed number of slots (defined by its type). In the slots, there are either the names of other chunks (denoting a connection between the two chunks) or primitive elements (i.e. symbols which do not refer to another chunk).

**Definition 6 (chunk store).** A chunk-store $\Sigma$ over a set of symbols $\mathcal{S}$ is a tuple $(\mathcal{C}, \mathcal{E}, \mathcal{T}, \mathcal{S}, \text{Slots}, \text{HasSlot}, \text{Isa})$, where $\mathcal{C}$ is a set of chunk identifiers and $\mathcal{E}$ a set of primitive elements both identified by unique names. The values of $\Sigma$ are defined by the set $\mathcal{V} = \mathcal{C} \cup \mathcal{E}$. The set $\mathcal{T}$ denotes the set of all chunk-type names and $\mathcal{S}$ the slot names. The function $\text{Slots} : \mathcal{C} \rightarrow 2^\mathcal{S}$ determines the allowed slots a chunk-type. The sets $\mathcal{C}, \mathcal{E}, \mathcal{T}$ and $\mathcal{S}$ are disjoint and $\mathcal{C} \cup \mathcal{E} \cup \mathcal{T} \cup \mathcal{S} \subseteq \mathcal{S}$.

$\text{HasSlot} : \mathcal{C} \times \mathcal{S} \rightarrow \mathcal{V}$ is a partial function determining the values of a chunk’s slots. $\text{Isa} : \mathcal{C} \rightarrow \mathcal{T}$ is a total function determining the type of a chunk, so each chunk has exactly one type. A chunk-store is type-consistent, iff the following two condition hold:

$$\forall c \in \mathcal{C}, s \in \mathcal{S}, v \in \mathcal{V} : (c, s, v) \in \text{HasSlot} \Rightarrow s \in \text{Slots}(\text{Isa}(c))$$

3.2 Buffer Systems

Buffer systems extend the definition of chunk stores by buffers. Each buffer can hold at most one chunk from its chunk store. This is modeled by the relation Holds in the following definition:

**Definition 7 (buffer system).** A buffer system is a tuple $(\mathcal{B}, \Sigma, \text{Holds})$, where $\mathcal{B}$ is a set of buffer names, $\Sigma = (\mathcal{C}, \mathcal{E}, \mathcal{T}, \mathcal{S}, \text{Slots}, \text{HasSlot}, \text{Isa})$ a type-consistent chunk-store and $\text{Holds} : \mathcal{B} \rightarrow \mathcal{C}$ a partial function that assigns every buffer at most one chunk that it holds. Buffers that do not appear in the $\text{Holds}$ relation are called empty. A buffer system is consistent, if $\Sigma$ is a type-consistent chunk-store. It is clean, if its chunk-store only contains chunks which appear in $\text{Holds}$.

3.3 The Operational Semantics of ACT-R

We define the operational semantics of ACT-R as a state transition system $(\mathcal{S}, \rightarrow)$. The state space $\mathcal{S}$ consists of states defined as follows:

**Definition 8 (ACT-R states).** An ACT-R state $s$ is a consistent buffer system $(\mathcal{B}, \Sigma, \text{Holds})^\mathcal{V}$ where $\mathcal{V}$ is a set of variable bindings.

The transition relation of our ACT-R semantics is based on a normal form, the variable normal form (VNF):
Definition 9 (variable normal form of an ACT-R rule). The variable normal form (VNF) of an ACT-R rule \((p \rightarrow LHS \Rightarrow RHS)\) is the term \((p \rightarrow LHS' \Rightarrow RHS')^{\mathcal{X}}\) where each symbol \(x_i\) appearing in LHS and RHS is replaced by a fresh variable \(X_i = x_i\). ACT-R variables of the form \(=x\) (not followed by the symbol \(\Rightarrow\)) are replaced by a fresh logical variable \(Y\).

Example 2 (VNF). The production rule \((p \rightarrow \mathbf{b} > s_1 \ s_2 \Rightarrow \mathbf{b} > s_1 \Rightarrow s_2)\) has the VNF \((p \rightarrow \mathbf{X}_1 > \mathbf{X}_2 \ \mathbf{X}_3 \ \mathbf{X}_4 \ \mathbf{X}_5 \Rightarrow \mathbf{X}_6 > \mathbf{X}_7 \ \mathbf{X}_8)^{\mathcal{Y}}\) with

\[
\mathcal{X} = \{X_1 = b, X_2 = s_1, X_3 = v_1, X_4 = s_2, X_5 = Y, X_6 = b, X_7 = s_1, X_8 = Y\}.
\]

In this example, logical variables are denoted with capital letters and constant symbols with small letters.

Definition 10 (operational semantics of ACT-R). For a production rule \(r\) with variable normal form \(P := (p \rightarrow LHS \Rightarrow RHS)^{\mathcal{Y}}\) where

\[
- \ LHS = \bigoplus_{i=1}^{n_r} \left( =b_i^L \Rightarrow \right. t_i \ \bigoplus_{j=1}^{n_r} \left( =s_i^L \ x_{i,j} \right) \bigoplus_{k=1}^{n_r} \left( =b_i^R \Rightarrow \right. \bigoplus_{l=1}^{n_r} \left( =s_i^R \ x_{k,l} \right) \bigoplus_{k=1}^{n_r} \left( =b_i^R \Rightarrow \right. \bigoplus_{l=1}^{n_r} \left( =s_i^R \ x_{k,l} \right)
\]

the transition \(\Rightarrow_r\) is defined as follows. Note that due to the variable normal form of the rule, all terms appearing in LHS and RHS do not contain constants.

If

- there is an instance of \(r^{\mathcal{X}}[X/Y]\) with variables \(X\) substituted by fresh variables \(Y\) and
- \(\forall (V \rightarrow \exists Y \ (\mathcal{X}[X/Y] \wedge C[X/Y] \wedge V[X/Y]))\) with
- \(C := \bigwedge_{i=1}^{n_r} \left( b_i^S \in B \wedge b_i^S \in B \wedge b_i^S \in B \wedge b_i^S \in B \right. t_i \ \bigwedge_{j=1}^{n_r} \left( =s_i^L \ x_{i,j} \right) \bigwedge_{k=1}^{n_r} \left( =s_i^R \ x_{k,l} \right) \bigwedge_{k=1}^{n_r} \left( =b_i^R \Rightarrow \right. \bigoplus_{l=1}^{n_r} \left( =s_i^R \ x_{k,l} \right) \bigoplus_{k=1}^{n_r} \left( =b_i^R \Rightarrow \right. \bigoplus_{l=1}^{n_r} \left( =s_i^R \ x_{k,l} \right)
\]

- \(V := \bigwedge_{i=1}^{n_r} \left( b_i^S = b_i^S \wedge t_i = t_i \wedge =s_i^L = s_i^L = s_i^L = s_i^L \right. x_{i,j} = v_{i,j})\),
- \(\Sigma := (\Sigma, \text{E}, \text{T}, \text{S}, \text{Slots}, \text{HasSlot}, \text{Isa})\)

then \(\langle B; \Sigma; \text{Holds}\rangle^{\mathcal{Y}} \Rightarrow_r \langle B; \Sigma'; \text{Holds}\rangle^{\mathcal{Y} \cup \mathcal{X} \cup \mathcal{V}}\) where

- \(\Sigma' := (\Sigma, \text{E}, \text{T}, \text{S}, \text{Slots}, \text{HasSlot}', \text{Isa})\)
- \(\text{HasSlot}' := (\text{HasSlot} - \left\{ (c, s, v) \in \text{HasSlot} \ | \ (b, c) \in \text{Holds} \wedge s = s_{k,l}^R \right\}) \cup \bigcup_{k=1}^{n_r} \left\{ (c_k, s_{k,l}^R, x_{k,l}') \right\} \wedge (b_k^R, c_k) \in \text{Holds}\)

The main part of the applicability condition is \(C\). Basically, the condition states that if the current state matches the left-hand side of a rule, the rule can be applied. When a rule is applied, the values in the state are overridden by the values on the right-hand side of the rule.
4 Translation of ACT-R rules to CHR

In this section, we define a translation function \( \text{chr}(\cdot) \) which translates ACT-R production rules and states to corresponding CHR rules and states. But first, we need another definition:

**Definition 11 (modification list).** For a production rule \( P := (p \ LHS \Rightarrow RHS) \) in VNF with

- \( LHS = \bigoplus_{i=1}^{n_L}(= b_i^L > \text{isa } t_i \bigoplus_{j=1}^{n_v^{L,i}}(s_{i,j}^L v_{i,j})) \)
- \( RHS = \bigoplus_{k=1}^{n_R}(= b_k^R > \bigoplus_{l=1}^{n_v^{R,k}}(s_{k,l}^R v_{k,l})) \)

the modification list \( \text{mod}(P) := \{ s_{i,j}^L | X \rightarrow \exists k, l : b_k^L = b_k^R \land s_{i,j}^L = s_{k,l}^R \} \).

The modification list contains all slots whose values appear on the r.h.s. of a rule, i.e. which are modified.

**Definition 12 (translation of production rules).** An ACT-R production rule \( P := (p \ LHS \Rightarrow RHS) \) in VNF where

- \( LHS = \bigoplus_{i=1}^{n_L}(= b_i^L > \text{isa } t_i \bigoplus_{j=1}^{n_v^{L,i}}(s_{i,j}^L v_{i,j})) \)
- \( RHS = \bigoplus_{k=1}^{n_R}(= b_k^R > \bigoplus_{l=1}^{n_v^{R,k}}(s_{k,l}^R v_{k,l})) \)

can be translated to a CHR rule of the form \( r \odot H_r \equiv G \cdot B \) as follows:

- \( H_k := \bigwedge_{i=1}^{n_L} \left( \text{buffer} (b_i^L, c_i^L) \land \text{chunk} (c_i^L, t_i) \right) \land \bigwedge_{s_{i,j} \notin \text{mod}(P)} \text{chunk_has_slot}(c_i^L, s_{i,j}^L, v_{i,j}) \)

  where the \( c_i^L \) are fresh variables
- \( B := \bigwedge_{k=1}^{n_R} \bigwedge_{l=1}^{n_v^{R,k}} \text{chunk_has_slot}(c_k^L, s_{k,l}^R, v_{k,l}) \) where \( c_k^L \) are fresh variables
- \( H_r := \bigwedge_{k=1}^{n_R} \bigwedge_{l=1}^{n_v^{R,k}} \text{chunk_has_slot}(c_k^L, s_{k,l}^R, x_{k,l}) \) where \( x_{k,l} \) are fresh variables
- \( G := \mathcal{X} \cup \mathcal{Y} \cup \mathcal{Z} \)
  - \( \mathcal{Y} := \{ c_k^L = c_i^L | i \in \{1, \ldots, n_L\}, k \in \{1, \ldots, n_R\}, \mathcal{X} \rightarrow b_i^L = b_k^R \} \)
  - \( \mathcal{Z} := \{ x_{k,l} = v_{i,j} | i \in \{1, \ldots, n_L\}, j \in \{1, \ldots, n_v^{L,i}\}, k \in \{1, \ldots, n_R\}, l \in \{1, \ldots, n_v^{R,k}\}, \mathcal{X} \rightarrow (b_i^L = b_k^R, s_{i,j}^L = s_{k,l}^R) \} \)

The translation of a production rule \( P \) according to this scheme is denoted as \( \text{chr}(P) \).

Note that \( H_r \) is constructed only by considering the right-hand side of the rule. Hence, there might be \( \text{chunk_has_slot} \) constraints in \( H_r \) with variables as values which are not bound in the guard \( G \). Those are the slots modified on the \( RHS \) but not appearing in \( LHS \) as mentioned in definition \( 11 \) and are not involved in the matching process of ACT-R rules (see definition \( 10 \)).

Informally, \( H_k \) contains all \( \text{buffer} \) and \( \text{chunk} \) constraints as well as all \( \text{chunk_has_slot} \) constraints of the slots which are not modified on the r.h.s. In contrast, \( H_r \) contains all \( \text{chunk_has_slot} \) constraints of the slots which appear on the r.h.s., i.e. which are modified. In the following example, the translation process is shown:
Example 3 (production rule translation). The production rule from example 1

\[ P := (p \text{ recognize-win} \Rightarrow =b> \text{ result win}) \]

can be translated to the following CHR rule. We replace variables by their ground values according to the variable bindings in X, Y and Z.

\[
\text{recognize-win} @ buffer(B^L, C^L) \land \text{chunk}(C^L, T) \land \text{chunk_has_slot}(C^L, S^L_1, V^L_1) \\
\land \text{chunk_has_slot}(C^R, S^R_1, X) \\
\iff B^L = b \land T = \text{game} \land \\
S^L_1 = \text{me} \land S^L_2 = \text{opponent} \land S^R_1 = \text{result} \land \\
V^L_1 = \text{rock} \land V^L_2 = \text{scissors} \land V^R_1 = \text{win} \land C^L = C^R | \\
\text{chunk_has_slot}(C^R, S^R_1, V^R_1).
\]

The removed head are simply the modified slot with a fresh variable X as value. This variable is singleton in this case, i.e. it does only appear once in the rule. This is because the result slot is only modified on RHS but never tested, hence the value does not influence the matching (and therefore, X does not appear in the set Z in the guard).

We now have defined how our subset of ACT-R production rules can be translated to CHR. As mentioned before, this translation formalism can be extended by requests and clearings: A request is a modification which modifies the complete chunk of a buffer, i.e. all constraints of that buffer are replaced. The replacement is achieved by a built-in function (in our case a Prolog predicate), which has a chunk description as input and then adds the corresponding constraints describing the result chunk. In the following definition, we present the translation of ACT-R states to CHR states.

Definition 13 (translation of states). An ACT-R state \( s := (B; \Sigma; \text{Holds})^V \) with \( \Sigma = (C, E, T, S, \text{Slots}, \text{HasSlot}, \text{Isa}) \) can be translated to the corresponding CHR state (denoted by \( \text{chr}(s) \)):

\[
\bigwedge_{(b,c) \in \text{Holds}} \text{buffer}(b,c) \land \bigwedge_{(c,t) \in \text{Isa}} \text{chunk_type}(c,t) \\
\land \bigwedge_{(c,s,v) \in \text{HasSlot}} \text{chunk_has_slot}(c,s,v) \land V.
\]

5 Soundness and Completeness

In this section, we prove soundness and completeness of our translation scheme from definition 12 and definition 13. I.e., we show that each transition of an ACT-R model in a certain state is also possible in the corresponding CHR program with the corresponding CHR state leading to the same results and vice versa. This is illustrated in figure 1. At first, we show that applicability is preserved by the translation and then extend this property to the soundness and completeness theorem.
Lemma 1 (applicability). If the production rule $P := (p \rightarrow LHS \Rightarrow RHS)^X$ in VNF where

$\begin{align*}
LHS &= \bigoplus_{i=1}^{n^L} (\langle b_i^L > \text{ isa } t_i^L \rangle \ \bigoplus_{j=1}^{n^L} (s_{i,j}^L, x_{i,j})) \\
RHS &= \bigoplus_{k=1}^{n^R} (\langle b_k^R > \ \bigoplus_{l=1}^{n^R} (s_{k,l}^R, x_{k,l}))
\end{align*}$

is applicable in ACT-R state $s$, then the corresponding CHR rule $chr(P)$ is applicable in state $chr(s)$.

Proof. Let $s = (B; \Sigma; \text{Has})$ be an ACT-R state with a chunk store $\Sigma := (C, E, T, S, \text{Slots}, \text{HasSlot}, \text{Isa})$. Rule $P$ is considered as a fresh rule instance. The rule $chr(P)$ has the form $r@H_k \Rightarrow G|B$.

$\Rightarrow$: If $r$ is applicable in $s$, then the applicability condition holds according to definition 10. I.e., the following conditions hold for $s$:

- $CET \models \forall (V \Rightarrow \exists Y (X \land C \land M))$
- $C := \bigland_{i=1}^{n^R} (b_i^S \in B \land b_i^S \text{ Holds } c_i^S \land c_i \in C \land c_i^S \text{ Isa } t_i^S, \ \
\land \bigland_{j=1}^{n^L} (c_i^S, s_{i,j}^S, v_i, j) \in \text{HasSlot})$
- $M := \bigland_{i=1}^{n^R} (b_i^L = b_i^S \land t_i^L = t_i^S \land \bigland_{j=1}^{n^L} (s_{i,j}^L = s_{i,j}^S, x_{i,j} = v_{i,j}))$

Hence, the corresponding CHR state $chr(s)$ contains at least the constraints $H \land V$ with

$H := \bigland_{i=1}^{n^R} (\text{buffer}(b_i^S, c_i^S) \land \text{chunk}(c_i^S, t_i^S) \ \
\land \bigland_{j=1}^{n^L} \text{chunk\_has\_slot}(c_i^S, s_{i,j}^S, v_{i,j}) \land R)$

by definition 13. According to definition 5, the applicability condition of the rule $chr(P)$ is

$CET \models \forall (D \Rightarrow \exists x (G \land H_k = H_k' \land H_r = H_r'))$

The rule is applicable, since every element in $H_r'$ has a matching partner defined by the matching $M_k := M \cup \{c_i^L = c_i^S | \forall i \in \{1, \ldots, n^L\}\}$ in $H$. For the constraints in $H_r'$, there are two cases:

**Case 1:** The constraint $c' := \text{chunk\_has\_slot}(c_i^L, s_{i,j}^L, x_{i,j}) \in H_r'^L$ is an instance of the constraint $c := \text{chunk\_has\_slot}(c_i^R, s_{k,l}^R, y_{k,l}) \in H_r'$ with $X \rightarrow b_i^L = b_k^R \land s_{i,j}^L = s_{k,l}^R$. In the guard, the fresh variables $c_i^L$ and $c_i^S$ are unified (which
is possible since \( c^R_k \) is unbound). Additionally, the fresh variables \( y_{k,l} \) are bound to \( x_{k,l} \), if \( X \to (b^L_{i,j} = b^R_{i,j} \land s^L_{i,j} = s^R_{i,j}) \). Since \( c^L_i = c^S_i \) and \( x_{i,j} = v_{i,j} \) is set in \( M_k \), the matching together with the variable bindings in the guard hold in state \( H \).

**Case 2:** For the other elements of \( H'_r \), i.e. all \( d := chunk\_has\_slot(c^R_k, s^L_{k,l}, y_{k,l}) \) \( \in H'_r \) with \( \overline{d} : X \to b^L_{i,j} = b^R_{i,j} \land s^L_{i,j} = s^R_{i,j} \), there are constraints \( buffer(b^L_{i,j}, c^L_i) \in H \) with \( b^L_{i,j} = b^R_{i,j} \), since on the r.h.s. of an ACT-R production rule, there are only buffers which appear on l.h.s. by definition \([1]\). In the guard, \( c^L_i \) is unified with \( c^R_k \) (which is possible, since \( c^R_k \) is unbound). The variables \( y_{k,l} \) remain unbound in the guard according to definition \([12]\). Due to \( c^L_i = c^S_i \) in \( M_k \), the match holds together with the variable bindings \( c^L_i = c^R_k \), since according to definition \([1]\) the slots \( s^R_{k,l} \) exist for the chunk type as defined on the l.h.s. of the rule and the HasSlot relation is left-total.

Hence, there is a matching for \( H'_k \) and \( H'_r \) respecting the guards and the rule is applicable in \( H \).

“\( \Rightarrow \)” Since \( chr(r) \) is applicable in \( chr(s) \Rightarrow H \), the state must have the form \( H = H_k \land H_r \land (H_k = H'_k) \land (H_r = H'_r) \land G \) with a corresponding matching \( M \). The corresponding ACT-R state \( s \) hence must have at least the following elements:

- for all constraints \( buffer(b^S_i, c^S_i) \in H_k \), there is a corresponding \( b^S_i \in B \) and \( c^S_i \in C \),
- for all constraints \( chunk\_has\_slot(c^S_i, s^S_{i,j}, v^S_{i,j}) \in H_k \), there is a \( (c^S_i, s^S_{i,j}, v^S_{i,j}) \in \text{HasSlot} \),
- for all constraints \( chunk\_has\_slot(c^S_i, s^S_{i,j}, v^S_{i,j}) \in H_r \), there is a \( (c^S_i, s^S_{i,j}, v^S_{i,j}) \in \text{HasSlot} \),
- The variable bindings \( V \cup X \) correspond to the matching \( M = \{ b^L_{i,j} = b^S_i, t^L_{i,j} = t^S_i, s^L_{i,j} = s^S_{i,j}, x_{i,j} = v_{i,j} \} \cup Y \cup Z \) with \( Y \) and \( Z \) according to definition \([12]\) due to the applicability of the rule.

Hence, the applicability condition of \([10]\) holds and the rule is applicable.

**Theorem 1 (soundness and completeness).** For a production rule of the form \( (p \to r) \text{ LHS} \Rightarrow \text{RHS} \) where

- \( \text{LHS} = \bigoplus_{i=1}^{n_v} (\text{isa } t_i \bigoplus_{j=1}^{n_b} (s^L_{i,j}, v^S_{i,j})) \)
- \( \text{RHS} = \bigoplus_{k=1}^{n_a} (\text{isa } b^R \bigoplus_{l=1}^{n_b} (s^R_{k,l}, v^S_{k,l})) \)

and two ACT-R states \( s \) and \( s' \) the transitions \( s \to_r s' \) and \( chr(s) \to_r chr(s') \) correspond to each other.

**Proof.** Let \( chr(P) = r @ H'_k \setminus H'_r \Leftrightarrow G \mid B \).

“\( \Rightarrow \)” According to lemma \([1]\) \( r \) is applicable in \( s \) iff \( chr(r) \) is applicable in \( chr(s) \). Let \( chr(s) \to_r s'' = (H_k \land C \land H_r = H'_k \land H_r = H'_r \land B \land G) \) (definition \([5]\)) It remains to show that the resulting state \( s'' \Leftrightarrow chr(s') \). Let \( s = (B; \Sigma; \text{Holds}) \).
be an ACT-R state with a chunk store \( \Sigma = (C, E, T, S, \text{Slots}, \text{HasSlot}, \text{Isa}) \) and 
\( s' = (B; \Sigma'; \text{Holds})^V \) with \( \Sigma' = (C, E, T, S, \text{Slots}, \text{HasSlot}', \text{Isa}) \). Then

\[
\text{HasSlot}' := \left( \text{HasSlot} - \left\{ (c, s, v) \in \text{HasSlot} \mid \mathcal{V} \cup X \rightarrow (b_k, c) \in \text{Holds} \land s = s^R_{k,l} \right\} \right) \\
\cup \bigcup_{k=1}^{n^R} \bigcup_{s=1}^{n^R} \left\{ (c_k, s^R_{k,l}, x'_{k,l}) \mid \mathcal{V} \cup X \rightarrow (b^R_k, c_k) \in \text{Holds} \right\}
\]

The resulting ACT-R state \( s' \) only modifies the HasSlot relation. We relate the above modifications of the state to their CHR counterparts according to definition 13:

- The removed elements in \( \text{chr}(s') \) are \( R := \bigwedge_{k=1}^{n^R} \bigwedge_{s=1}^{n^R} \text{chunk-has_slot}(c^R_k, s^R_{k,l}, \cdot) \).
- The added elements in \( \text{chr}(s') \) are \( A := \bigwedge_{k=1}^{n^R} \bigwedge_{s=1}^{n^R} \text{chunk-has_slot}(c^R_k, s^R_{k,l}, x'_{k,l}) \).

The elements \( R \) correspond to \( H_r \) of the state. All other constraints from \( \text{chr}(s) \) are kept in \( s'' \) which also corresponds to \( \text{chr}(s') \). The added constraints \( A \) correspond to the added constraints in \( s'' \) which are

\[ B = \bigwedge_{k=1}^{n^R} \bigwedge_{s=1}^{n^R} \text{chunk-has_slot}(c^R_k, s^R_{k,l}, x'_{k,l}) \]

The added built-ins are \( G = \mathcal{X} \cup \mathcal{Y} \cup \mathcal{Z} \) according to definition 12. The built-ins in \( \mathcal{X} \) are also added to \( \mathcal{Y} \) by definition 10 and hence are part of \( \text{chr}(s') \). The matching \( M \) which is added to \( s'' \) in form of \( H'_k = H_k \) and \( H'_z = H_z \) is also added to \( s' \). \( \mathcal{Y} \) and \( \mathcal{Z} \) only bind fresh variables to variables in \( \text{chr}(s') \). The result is equivalent to the translation \( \text{chr}(s') \), where the bound variables are reused in the first place according to definition 13.

"\( \Rightarrow \)" Let \( \text{chr}(s) \Rightarrow_r \text{chr}(s') \) and \( s \Rightarrow_r s'' \). According to lemma 1, \( \text{chr}(r) \) is applicable in \( \text{chr}(s) \) iff \( r \) is applicable in \( s \). It remains to show that the resulting state \( s'' = s' \). We define the following variables:

\[
s = (B; \Sigma; \text{Holds})^V, \quad \Sigma = (C, E, T, S, \text{Slots}, \text{HasSlot}, \text{Isa})
\]

\[
s' = (B; \Sigma'; \text{Holds})^V, \quad \Sigma' = (C, E, T, S, \text{Slots}, \text{HasSlot}', \text{Isa})
\]

\[
s'' = (B; \Sigma''; \text{Holds})^{V''}, \quad \Sigma'' = (C, E, T, S, \text{Slots}, \text{HasSlot}'', \text{Isa})
\]

According to definition 10

\[
\text{HasSlot}'' = \left( \text{HasSlot} - \left\{ (c, s, v) \in \text{HasSlot} \mid (b, c) \in \text{Holds} \land s = s^R_{k,l} \right\} \right)
\\
\cup \bigcup_{k=1}^{n^R} \bigcup_{s=1}^{n^R} \left\{ (c_k, s^R_{k,l}, x'_{k,l}) \right\}
\]

By definition 5, \( \text{chr}(s') = (H_k \land C \land H_k = H'_k \land H_r = H'_r \land B \land G) \). with \( H'_k, H'_z, G \) and \( B \) defined as in definition 12. This leads to the following ACT-R state \( s' \):

- The constraints in \( C \) are kept. The corresponding in \( s' \) are copied from \( s \).
- The \( buffer, chunk \) and \( chunk-has_slot \) constraints in \( H_k \) remain in \( \text{chr}(s') \).
- The corresponding elements in \( s' \) are copied from \( s \).
– The chunk_has_slot\( (c, S_{i,j}, v_{i,j}) \) constraints in \( H_r \) are removed from \( chr(s') \).

The corresponding elements in \( s \) are removed from \( s' \) (in particular from HasSlot').

– \( V'' = \varnothing \).

This state \( s' \) is equivalent to \( s'' \), since the same elements are kept, removed and added (analogously to the other direction of the proof). The variable bindings are equivalent as argued above.

6 Example: Rock, Paper, Scissors

Example 4 (rock, paper, scissors in ACT-R). Consider the game rock, paper, scissors. To represent this game in ACT-R, we introduce – analogously to example 1 – a chunk type game with slots me representing the model’s move, opponent representing the opponent’s move and result storing if a round has been won or lost by the model or if it is a draw round. The moves are symbols rock, paper and scissors, the possible results are win, fail, draw.

We assume that the game chunk of a round is initialized with an empty me slot (represented by the symbol nil) and a move in the opponent slot. Then our model has three choices to play. Each choice is modeled by a production rule:

\[
\begin{align*}
(p \text{ play-move}) \\
=b> \text{isa game me nil} \Rightarrow =b> \text{me move}
\end{align*}
\]

where \( move \in \{\text{rock, paper, scissors}\} \). Hence, the model is non-deterministic, since more than one rule is applicable in the initial state according to the operational semantics (definition 10).

For each possible game situation (combinations of rock, paper, scissors in the me and opponent slots), there is a rule recognizing a win, failure or a draw. For instance, consider a variant of the rule from example 1 recognizing a win situation:

\[
\begin{align*}
(p \text{ recognize-win}) \\
=b> \text{isa game me rock opponent scissors result nil} \Rightarrow =b> \text{result win}
\end{align*}
\]

The play-move rules can be translated analogously to the following CHR rule for play-rock (not in HNF for better readability) using 12:

\[
\text{play-rock @ buffer}(b, C) \land \text{chunk}(C, \text{game}) \setminus \text{chunk_has_slot}(C, \text{me, nil}) \\
\Rightarrow \text{chunk_has_slot}(C, \text{me, rock}).
\]

The recognize-win rule (and analogous rules) can be translated as in example 1.

We now consider some transitions which are possible in a given initial ACT-R state \( s_0 := \{(b), \Sigma, \{(b, c)\}\} \) with

\[
\Sigma = \{(c), \{\text{rock, paper, scissors, win, draw, fail}, \text{game}\}, \{\text{me, opponent, result}\}, \{(\text{game}, \{\text{me, opponent, result}\})\}, \{(\text{c, me, nil}), (\text{c, opponent, rock}), (\text{c, result, nil})\}, \{(\text{c, game})\})
\]
I.e., \( s_0 \) is a state for an ACT-R system with one buffer \( b \) which holds a chunk \( c \) of type \textit{game} where the opponent played \textit{rock} and the other values are not initialized. The following transition is possible (we only regard the \texttt{HasSlot} relation):

\[
\{(c, \text{me}, \text{nil}), (c, \text{opponent}, \text{rock}), (c, \text{result}, \text{nil})\} \xrightarrow{\text{play-rock}} \{(c, \text{me}, \text{rock}), (c, \text{opponent}, \text{rock}), (c, \text{result}, \text{nil})\}
\]

\[
\xrightarrow{\text{recognize-draw}} \{(c, \text{me}, \text{rock}), (c, \text{opponent}, \text{rock}), (c, \text{result}, \text{draw})\}
\]

The other two possible transitions lead to the following final states:

- \( \{(c, \text{me}, \text{paper}), (c, \text{opponent}, \text{rock}), (c, \text{result}, \text{win})\} \) (using \texttt{play-paper} and \texttt{recognize-win})
- \( \{(c, \text{me}, \text{scissors}), (c, \text{opponent}, \text{rock}), (c, \text{result}, \text{fail})\} \) (using \texttt{play-scissors} and \texttt{recognize-fail})

The corresponding CHR program has the following transitions for the initial state \( \texttt{chr}(s_0) \) from above (we only regard the \texttt{chunk_has_slot} constraints abbreviated by \texttt{chs}):

\[
\texttt{chs}(c, \text{me}, \text{nil}) \land \texttt{chs}(c, \text{opponent}, \text{rock}) \land \texttt{chs}(c, \text{result}, \text{nil}) \xrightarrow{\texttt{play-rock}} \texttt{chs}(c, \text{me}, \text{rock}) \land \texttt{chs}(c, \text{opponent}, \text{rock}) \land \texttt{chs}(c, \text{result}, \text{nil})
\]

\[
\xrightarrow{\texttt{recognize-draw}} \texttt{chs}(c, \text{me}, \text{rock}) \land \texttt{chs}(c, \text{opponent}, \text{rock}) \land \texttt{chs}(c, \text{result}, \text{draw})
\]

The other two transitions can be derived analogously.

7 Conclusion

In this paper, we have presented a formalization of the core of the production rule system ACT-R including an abstract operational semantics. Furthermore, we have shown a formal translation of ACT-R production rules to CHR. The translation is sound and complete.

The formalization of ACT-R is based on prior work. In [6] we have presented an informal description of the translation of ACT-R production rules to CHR rules. This informal translation has been implemented in a compiler transforming ACT-R models to CHR programs. Our implementation is modular and exchangeable in its core features as we have shown in [7] by exchanging the central part of the conflict resolution with four different methods. Although the implementation is very practical and covers a lot of practical details of the ACT-R implementations, it is not directly usable for analysis.

Our formalization of the translation process in this paper is very near to the practical implementation as it uses the same translation schemes for chunk stores, buffer systems and consequently states. Even the rules are a simplified version of our practical translation from [7]. However, it abstracts from practical aspects like time or conflict resolution. This is justifiable, since for confluence analysis, this kind of non-determinism in the operational semantics is useful. Additionally,
as shown in section 6, the general computation process is reproduced closely by our semantics. Furthermore, due to the soundness and completeness of our translation, confluence analysis tools from CHR can be used on our models.

Hence, the contributions of this paper are

– an abstract operational semantics of ACT-R which is – to the best of our knowledge – the first formal representation of ACT-R’s behavior,
– a formal description of our translation process (since in 6 a more technical description has been chosen),
– a soundness and completeness result of the abstract translation.

For the future, we want to extend our semantics such that it covers the more technical aspects of the ACT-R production rule system like time and conflict resolution. We then want to investigate how this refined semantics is related to our abstract operational semantics from this paper. The example in section 6 is non-deterministic and not confluent since the result state varies depending on the move of the model. The confluence test described in 4, pp. 102 sqq., reveals the same result. To overcome non-determinism, ACT-R uses a conflict resolution strategy. In 7 we have analyzed several conflict resolution strategies using the same example as in section 6. A confluence test might be useful to reveal rules where the use of conflict resolution is undesired. For the future, we want to investigate how the CHR analysis tools perform for our ACT-R semantics and how they might support modelers in testing their models for undesired behavior, since the application of the confluence test on our example is promising.

Additionally, it could be interesting to use the CHR completion algorithm 11 to repair ACT-R models that are not confluent.

References

A Visualization Tool for Constraint Handling Rules

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Abstract. Constraint Handling Rules (CHR) is currently used as a general purpose language. This results in having complex programs with CHR. Nevertheless, CHR is still lacking on visualization tools. Such tools are useful for beginners to the language as well as programmers of difficult algorithms. This paper continues upon the efforts made to have a generic visualization platform for CHR using source-to-source transformation. It also provides a new visualization feature that enables viewing all the possible solutions of a CHR program instead of the don’t care nondeterminism used in most CHR implementations.

Keywords: Constraint Handling Rules, Algorithm Visualization, Source-to-Source Transformation

1 Introduction

Constraint Handling Rules (CHR) \(^1\) is a committed-choice rule-based language with multi-headed rules. It rewrites constraints until they are solved. (CHR) has developed from a language for writing constraint solvers into a general purpose language. Different classes and types of algorithms are currently implemented using CHR.

So far, visually tracing the different algorithms implemented in CHR was not possible. Debugging the algorithms implemented through the different CHR programs was not visualized. Such visual tools are important for any programming language. The lack of such tools makes it harder for programmers to trace complex algorithms that could be implemented with CHR. Previous work was directed towards the visualization of the execution of the different CHR rules in a program. Visual CHR \(^2\) was able to visualize the execution of CHR programs. However, it was directed towards the Java implementation of CHR; JCHR \(^3\). In addition, the compiler of JCHR had to be modified to add the required features. In \(^4\) a new tool was provided. It was able to visualize the execution of CHR programs implemented through Prolog. The tool used source-to-source transformation to eliminate the need to modify the compiler.

Although the tool provided through \(^4\) was able to add some visualization features to CHR, it lacked generality. It was only able to visualize the execution
of the different rules in a step-by-step manner. In addition, it was able to visualize CHR constraints as objects. However, the choice of the objects was limited and the specification of the parameters of the different objects was very rigid. Although the tool could be extended to animate the execution of different algorithms, the need of having static inputs remained due to the inflexibility of the provided tracer. The attempts provided through [5] and [6] also had the problem of being tailored to some specific algorithms.

Thus the tool presented through this paper aims at providing a more general CHR visualization platform. In order to have a flexible tracer, it was decided to use an already existing visualization tool. Such tools usually provide a wide range of objects and sometimes actions as well. As a proof of concept, we used Jawaa [7] throughout the paper. The annotation tool is available through: https://sourceforge.net/projects/chrvisualizationtool

In addition to introducing a more generic CHR algorithm visualization system, the tool also has a module that allows the user to visualize the exhaustive execution of any CHR program thus forcing the program to produce all the possible solutions. This allows the user to trace the flow of a CHR program using a different semantics than the refined operational semantics [8] embedded in SWI-Prolog. The output of the visualization is a tree showing the different paths of the algorithm’s solutions. It is also linked to the visualization tool as shown in details in Section 6.

The paper is organized as follows: Section 2 introduces CHR. Section 3 introduces the details of the annotation module. The details of the transformation approach are presented in Section 4. Section 5 shows an example of the visualization of an algorithm implemented through CHR. Section 6 shows how it was possible to transform CHR programs to produce all the possible solutions instead of only one. Finally, we conclude with a summary and directions for future work.

2 Constraint Handling Rules

This section introduces CHR through an example of a program that solves the nqueens problem.

There are two types of constraints: CHR constraints introduced by the user and built-in constraints [9]. A CHR program consists of a set of rules. Each rule has a head, a body and an optional guard. The head of any CHR rule consists of a conjunction of CHR constraints. The guard of a rule is used to set a condition for applying the rule. The guard can thus only contain built-in constraints. The body, on the other hand, can contain both CHR and built-in constraints [9].

In general there are three types of CHR rules: simpagation, propagation and simplification rules. A simpagation rule has the form:

\[ \text{optional_rule_name} \rightarrow H_K \setminus H_R \Leftrightarrow G \mid B. \]

As seen, there are two types of head constraints. \( H_K \) is the conjunction of CHR constraint(s) that are kept once the rule is executed. On the other hand, \( H_R \) contains CHR constraint(s) that are removed once the rule is executed. \( G \) is the optional guard that has to be satisfied to execute the rule. \( B \) is the body of
the rule. The constraints in \( B \) are added to the constraint store once the rule is executed.

A simplification rule is a simpagation rule with empty \( H \). This implies that the head constraint(s) are removed on executing the rule. A simplification rule has the following form:

\[
\textit{optional rule name} @ H_R \iff G \mid B.
\]

Finally, a propagation rule is a simpagation rule with empty \( H_R \). Thus, on executing a propagation rule, its body constraints are added to the constraint store without removing anything. Its format is:

\[
\textit{optional rule name} @ H_K \implies G \mid B.
\]

The following program implements an algorithm that solves the \( n \) queens problem \(^1\) which tries to place \( n \) queens in an \( n \) by \( n \) board such that no queen can attack another. Two queens can attack each other vertically, horizontally or diagonally.

\[
\begin{align*}
\text{initial} @ \text{solve}(N) & \iff \text{generate}(1,N,\text{List}), \text{queens}(N,\text{List}), \text{labelq}. \\
\text{add1} @ \text{queens}(N,\text{Dom}) & \iff N \gt 0 \mid N_1 \text{ is } N-1, \text{in}(N, \text{Dom}), \text{queens}(N_1,\text{Dom}). \\
\text{add2} @ \text{queens}(0,\text{Dom}) & \iff \text{true}. \\
\text{reduce} @ \text{in}(N1, [P]) \setminus \text{in}(N2, \text{Dom}) & \iff P_1 \text{ is } P-(N1-N2), P_2 \text{ is } P+(N1-N2), \text{delete}(\text{Dom},P,D1), \text{delete}(D1,P1,D2), \text{delete}(D2,P2,D3), \text{Dom} \iff D3 \mid D3 \iff \text{true}, \text{in}(N2, D3). \\
\text{label} @ \text{labelq} \setminus \text{in}(N, \text{Dom}) & \iff \text{Dom}=\text{[\_\_\_\_]} \mid \text{member}(P,\text{Dom}), \text{in}(N, [P]).
\end{align*}
\]

The model of the problem uses \( n \) variables each represented using the \textit{queens/2} constraint. The value of every variable determines the row number. The index, on the other hand, determines the column number. For example if the value of the second queen is three, this means that the queen in the second column is placed in the third row. The domain of any queen is initialized to be from 1 to \( N \) using the predicate \textit{generate/3}.

As seen from the program the rule \textit{initial} is used to initialize the solving process. The query the user enters contains the constraint \textit{solve/1}. Thus the, initially empty, constraint store is filled with the constraint in the query at the beginning. Afterwards, the rule \textit{initial} removes the constraint \textit{solve/1} since it is a simplification rule. It adds the two constraints \textit{queens} and \textit{labelq} which enable finding a solution. As seen from the rule, the second argument of the \textit{queens} constraint is set to be a list containing all the numbers from 1 till \( N \).

The two rules \textit{add1} and \textit{add2} are used to initialize the domains of all of the queens of the board using the previously computed list. The domain of every queen is represented using the \textit{in/2} constraint.

The rule \textit{reduce} is used to prune the domains of the different queens. In order to execute the rule, the location of a specific queen has to be determined. This is represented by having a domain list with one element only. The rule removes from the domain of another queen any value that could lead to an attack. This

\(^1\) The program is adapted from WebCHR \( \text{http://chr.informatik.uni-ulm.de/~webchr/} \)
ensures that whenever a location is chosen for this queen, it does not threaten the already labeled queen.

Finally the rule label is used to search through the domains whenever domain pruning is not enough.

3 Annotation to Visualize CHR Algorithms

According to [10], algorithm animation represents the different states of the algorithm as pictures. The animation differs according to the interaction between such states. To overcome the problems faced in [4], it was decided to use a generic visualization tracer that is able to produce different visual objects and possibly actions. However, in order to produce the needed visual states of the executed algorithms, the CHR programs had to be linked to the visual tracer. In other words, the transformation has to find a way to instruct the tracer to produce the needed visualization. The problem was how to achieve this goal while keeping a generic platform that is not tailored according to the algorithm type. The new system thus introduces a new module called the “Annotation Module”. Such module is needed to link between the different CHR constraints and the Jawaa objects/commands.

The idea is similar to the “interesting events” that BALSA[11] uses. In BALSA, the animator and the algorithm designer have to identify interesting events that could change the visualized images. Zeus [12], also uses the notion of interesting events and annotates the algorithm with markers that identify the basic operations. The preprocessor reads the event specifications and generates the needed definitions.

The Annotation module the system offers is however simpler and more flexible. It does not require having an animator who is different than the algorithm programmer. In the system, an interesting event is basically defined as the addition of CHR constraint(s) that leads to a change in the state of the algorithm and thus a change in the visualized data structure. For example, in sorting algorithms, every time an element in the list is moved to a new position, the list changes and thus the visualized list should change as well. This is an example of a CHR rule than can sort a list:

\[
\text{sort } @ \text{list}(I_1,V_1),\text{list}(I_2,V_2)<=>I_1<I_2,V_1>V_2|\text{list}(I_2,V_1),\text{list}(I_1,V_2)\]

As seen from the program, an element in the list could be represented by a constraint \text{list}/2 with the index and the value. Thus adding or changing the \text{list} constraint is basically the interesting event in this sorting algorithm.

For the \text{nqueens} algorithm presented in Section 2, two interesting events could be identified. The first one is through the \text{solve}/1 constraint which can be associated with initializing a visual \(N \times N\) board. The second interesting event can be associated with the \text{in}/2 constraint that represents the domains of the queens. Thus every time the domain changes a new visual state is produced and shown.

The annotation module provides its users with an interface through which they can choose to link constraint(s) with object(s) and/or action(s). In order to
have a dynamic system, the tool is automatically populated through a text file that contains the available objects and actions and their corresponding parameters. For example, the line `circle#name#x#y#width#color#bkgrd`, adds the object `circle` as one of the available objects to the user. The `circle` object requires the parameters `name`, `x`, `y`, `width`, `color` and `bkgrd`. The user can then enter the name of the constraint and the corresponding annotation as shown in Figure 1.

In this example the user associates the `in` constraint with the Jawaa object “Node”. As seen from the figure, the keyword “valueOf” can be used to assign to any of the needed parameters a value that is calculated through the argument(s) of the constraint. In addition, the user can specify a condition for activating the annotation. This condition has to hold in order for the corresponding visual step to be produced. In the shown example, the annotation is only valid if the length of the `List` is greater than one. In the given example, the name of the Jawaa node is “node” concatenated with the value of the first argument. Thus for the constraint `in(1,[1,2,3,4])`, the corresponding node has the name `node1`. The y-coordinate is also derived from the value of the first parameter but this time it is multiplied by `30`. The text inside the node also uses the value of the second parameter of the constraint.
Once the annotation is saved, a text file with all of the entered annotations is produced. For example, the previous annotation produces a text file containing:
```
in(Q, List) ==> node#length(valueOf(arg1), Len), Len > 1
#name=nodevalueOf(arg0)#x=160#y=valueOf(arg0)*30#width=90#height=30
#n=1#data=qvalueOf(arg0):valueOf(arg1)#color=black#bkgrd=green
#textcolor=black#type=RECT
```

In addition to the basic annotations shown in Figure 1, it is sometimes useful for users to add more advanced annotations. The rest of the section goes through the more complex annotations provided through the system.

The first possibility is that users can combine multiple constraints in one annotation. This means that the interesting event is not associated with having only one constraint in the store, it is rather having a number of constraints together in the constraint store. For example, users can specify that the color red and the color yellow result in the color orange. With such facility, users can specify all of the rules of the color mixing algorithm by using such annotations.

In addition, users can also specify more than one association to the same constraint. Such associations could be objects and/or actions. For example the user can specify that a `solve(N)` constraint is associated with \( N^2 \) rectangles. This way whenever the constraint `solve` is added to the store, \( N^2 \) rectangles are drawn. This results in drawing the initial board for the nqueens algorithm.

In addition, users can choose to annotate CHR rules instead of only having annotations to constraints. This results in adding Jawaa objects and/or actions whenever a specific rule is executed. In this case the interesting event is the execution of the rule as opposed to adding a constraint to the store. Thus whenever a rule is annotated this way, a new step in the visual trace is added on executing the rule.

### 4 Transformation Approach

The transformation approach is similar to the one presented in [13] and [4]. Both approaches present the CHR program using a set of constraints that encode the constituents of the CHR rule. For example `head(initial, solve(N), remove)` encodes the information that `solve(N)` is one of the head constraints of the rule `initial` and that this constraint is removed on executing the rule since it is a simplification rule.

The CHR program is thus first parsed to extract the needed information about the rules. The information is represented in the needed “relational normal” [13] form for the transformer to use. The transformation mainly aims at interfacing the CHR programs with the annotations that the user had already entered producing the needed visual states. For an input CHR program with the form:

```
:- chr_constraint cons/2.
initialrule @ cons(X1, X2) \ cons(X3, X4) <= somecheck | cons(X5, X6).
```

The output has the form:
The transformation thus adds for every constraint `constraint/N` a rule of the form:

```prolog
constraint(X1,X2,...,Xn) => check(status, false) | communicate_constraint(constraint(X1,X2,...,Xn)).
```

This extra rule makes sure that every time a new `constraint` is added to the constraint store, it is communicated to the tracer. In the case where the user had specified this `constraint` to be an interesting event (i.e. entered an annotation for it), it is used to produce the corresponding object(s)/action(s).

The original rules in the program are also modified to communicate the needed constraints. As seen from the previous rules, any new constraint added to the store is automatically communicated to the tracer. Thus once the body constraints are added to the store, they are automatically communicated to the tracer. Hence the rules can affect the visualization only through the head constraints. To be more specific, only the head constraints removed from the store can affect the resulting visualization since the corresponding object(s) might need to be removed as well from the visual trace. Thus the transformer can instruct the new rules to communicate the head constraints.

Whenever the user specifies a compound annotation associating the constraints `cons1(X11,...,X1n), cons2(X21,...,X2n),...,consn(Xn1,...,Xnn)` with `consnew(Xnew1,...,XnewN)`, a new CHR rule is added to the program. The rule has the form:

```prolog
cons1(X11,...,X1n), cons2(X21,...,X2n),...,consn(Xn1,...,Xnn) => consnew(Xnew1,...,XnewN)
```

As seen, the default case is to keep the constraints but the transformer can be instructed to produce a simplification rule instead.

The last case is the rule annotation. The problem with this case is that the CHR constraints in the body should be neglected since the whole rule is being annotated. Thus, even if the constraints were determined by the user to be interesting events, they have to be ignored since the execution of the rule includes them and the rule itself was annotated as an interesting event. Hence, to avoid having problems with this case, a generic status is used throughout the new program. In the transformed program, any rule annotated by the user changes the status to true at execution. As noticed from the first example in the section, all of the new rules added to the transformed program check that

---

2 The tracer is able to handle the problem of having multiple Jawaa objects with the same name by removing the old object having the same name before adding the new one. This is possible even if the removed head constraint was not communicated.
the status is set to false before communicating the corresponding constraint to the tracer. Thus whenever an annotated rule is executed, and on adding the constraints in the body to the store, those newly added rules are not triggered since their guards are not satisfied. This ensures that only the event associated with the rule is considered.

For the program:

\[
\text{hk} \implies \text{check} \left( \text{status}, \text{false} \right) \mid \text{communicate} \left( \text{hk} \right).
\]
\[
\text{hr} \implies \text{check} \left( \text{status}, \text{false} \right) \mid \text{communicate} \left( \text{hr} \right).
\]
\[
\text{annotationr1} \iff \text{communicate} \left( \text{annotationr1} \right).
\]
\[
\text{r1} \equiv \text{hk} \mid \text{hr} \iff \text{set} \left( \text{status}, \text{true} \right) \mid \text{body}, \text{annotationr1}, \text{set} \left( \text{status}, \text{false} \right).
\]

Only the rules that were annotated set the status to be true. This scheme is thus the most generic scheme that is currently being used. Consequently, in general even if no rules were annotated, the status check is always included.

5 Jawaa Example

This section shows how visual features were embedded to the queens algorithm introduced in Section 2. The full list of annotations is shown in Figures 2 and 3.

As seen from the annotation, the user can use the reserved keyword prologValue to compute a value (named “X”) using Prolog. The visual board is initialized through specifying that the solve constraint is an interesting event.
The `solve` constraint is annotated with 16 “rectangle” objects. Thus whenever the `solve` constraint is added to the store the 16 rectangles are visualized showing the initial board. This annotation is hence used with the 4-queens problem. The board consists of 16 adjacent rectangles each with width 30 and height 30. The `in/2` constraint has two different annotations. The first one is activated whenever the length of the list is equal to 1. This is the case where the queen is labeled to be placed in a specific position on the board. In this case, the x-coordinate of the Jawaa node is calculated as the index multiplied by the width of the cell which is 30. The y-coordinate is calculated through the only value in the domain i.e. the assigned value. It is also multiplied by 30. The second annotation is activated whenever the length of the domain list is greater than 1. In this case the queen is not placed in any position in the board since there are multiple possibilities. It is visualized as a “Node” outside the board and the domain is written on it. Figure 4 shows the visual steps produced for the query `solve(4)` until a solution is found.

For the sorting program given in Section 3, in addition to annotating the `list` constraint with a Jawaa “Node”, the user can also associate the rule `sort` with an annotation that highlights the bars before moving them to the correct positions. The visualization of the query `list(1,7),lit(2,6)` is shown in Figure 5 where the initial numbers 7 and 6 are shown with green bars at the first step. Afterwards, since the two numbers are not sorted with respect to each other, the rule `sort` is fired activating its corresponding annotations hence highlighting the two bars as shown in the second step. As a final step, the user is shown the two numbers after being sorted.

### 6 Visualizing Different Semantics

Although SWI-Prolog implements the refined operational semantics [8] for CHR, there are different proposed and defined CHR operational semantics. Based on the conflict resolution approach presented in [9], it is possible to convert a program running with a different operational semantics into the refined operational semantics used in SWI-Prolog. The abstract operational semantics of CHR [9] is non-deterministic. At any point, if several rules are applicable, one of them is randomly chosen. The application of a rule, however, cannot be undone since it is committed choice. In addition, the goal constituents are randomly chosen for processing. The refined operational semantics [8], on the other hand, chooses
Fig. 4: Visualizing the execution of the $n$queens algorithm for 4 queens.

A top-bottom approach for deciding on the applicable rule i.e. the first applicable rule is always chosen. In addition, the constraints are processed from left to right.
Fig. 5: Visualization of sorting algorithm.

Nondeterminism is especially interesting when the CHR program is non-confluent. Confluence [14] is a property of CHR which ensures the same final result no matter which applicable rule was chosen at any point of the execution.

The tool includes a module that is able to embed some of the non-determinism properties into any CHR solver. The newly generated solvers are able to choose, at any point of the execution, any of the applicable rules producing all the possible solutions.

6.1 Transformation Approach

This section discusses how any CHR program is transformed into a new one that is able to generate all the possible solutions of a non-confluent CHR program instead of using the refined operational semantics that generates only one solution. The transformation approach is based on the approaches presented in [9] and [15]. The main difference is that the new solver communicates some of the information to the visual tracer to be able to produce the needed visualization.

The transformed program starts each step by collecting the set of applicable rules with its corresponding head constraints. After the candidate list is built, the solver chooses one of the rules randomly using the built-in predicate select/3.

The newly transformed program is thus a CHR [16] solver. For example a rule of the form:

\[ r1 @ Hk \setminus Hr \iff \text{Guard} \mid \text{Body}. \]

generates two rules in the transformed program. The first generated rule is used to populate the candidate list. It is a propagation rule of the form:

\[ Hk, Hr \iff \text{Guard} \mid \text{cand}([(r1, [Hk, Hr]])]. \]

The second rule is fired whenever this rule is chosen from the candidate list. It has the following form:

\[ Hk\setminus\text{fire}((r1, [Hk, Hr])), Hr \iff \text{Guard} \mid \text{communicate_heads_kept}(Hk), \text{communicate_heads_removed}(Hr), \text{communicate_body}(body), \text{Body}. \]

In addition, the new program contains the following two rules:

\[ \text{cand}(L1), \text{cand}(L2) \iff \text{append}(L1, L2, L3) \mid \text{cand}(L3). \]

\[ \text{cand}([H|T]), \text{fire} \iff \text{select}(\text{Mem}, [H|T], \text{NList}), \text{fire}(\text{Mem}), \text{cand}(\text{NList}), \text{fire}. \]

The first rule ensures that the candidate list is correctly populated and incremented. The second rule, on the other hand, selects one of the elements of the candidate list at each step.

For example the program:

\[ : -\text{chr_constraint sphere/2}. \]

\[ r1 @ \text{sphere}(X,\text{red}) \iff \text{sphere}(X,\text{blue}). \]

\[ r2 @ \text{sphere}(X,\text{red}) \iff \text{sphere}(X,\text{green}). \]

is transformed into
:-chr_constraint sphere/2, fire/1, cand/1, fire/0.
r1 @ sphere(X,red) ==> cand([r1,[sphere(X,red)]]).
r2 @ sphere(X,red) ==> cand([r2,[sphere(X,red)]]).
cand(L1),cand(L2) <=> append(L1,L2,L3), cand(L3).
cand([H|T]),fire<=> select(Mem,[H|T],NList), call(fire(Mem)),
cand(NList),fire.
r1 @ fire([r1,[sphere(X,red)]]),sphere(X,red) <=>
  communicate_head_removed([sphere(X,red)]),
  communicate_body([sphere(X,blue)]), sphere(X,blue).
r2 @ fire([r2,[sphere(X,red)]]),sphere(X,red) <=>
  communicate_head_removed([sphere(X,red)]),
  communicate_body([sphere(X,green)]), sphere(X,green).

6.2 Visualization

With the refined operational semantics, the query sphere(a,red) results in executing r1 adding to the store the new constraint sphere(a,blue).

When visualizing the execution of the solver with this query, using the tool presented in [4], the result is shown in Figure (a). The CHR constraints remaining in the constraint store are shown in white and those removed are shown in red.

(a) Visualizing the execution of the solver.

(b) Showing all possible paths.

Fig. 6: Different options for visualizing the execution.

On the other hand, the transformed program is able to generate the visual tree shown in Figure (b). Since there were two applicable rules, the output tree accounts for both cases by the different paths. Through SWI-Prolog the user can trigger this behavior using the “;” sign to search for more solutions.

Given the solver:

rule1 @ sphere(X,red)<=>sphere(X,blue).
rule2 @ sphere(X,blue) <=> sphere(X,green).

The steps taken to execute the query sphere(b,blue), sphere(a,red) with the solver:

- First Solution
  1. In the first solution rule2 is fired replacing the constraint sphere(b,blue) by sphere(b,green).
  2. rule1 is then fired removing the constraint sphere(a,red) and adding the constraint sphere(a,blue).
  3. Finally, sphere(a,blue) triggers rule2 replacing it by sphere(a,green).

- Second solution
  1. Backtracking is triggered through the semicolon (;). We thus go back to the root and choose to apply rule1 for sphere(a,red) producing the sphere(a,blue).
  2. Afterwards, rule2 is executed to replace sphere(a,blue) by sphere(a,green).
  3. Finally, rule2 is fired replacing sphere(b,blue) by sphere(b,green).

- Third Solution
  1. This time when the user backtracks, execution goes back to the second level, applying rule2 to replace sphere(b,blue) by sphere(b,green).
  2. Afterwards, rule2 replaces sphere(a,blue) by sphere(a,green).

As seen from the tree in Figure 7, the constraint store in the final states is always the same containing sphere(a,green), sphere(b,green). However, the paths taken are different. Once the user enters a query, the visual trees are automatically shown.

In addition whenever the user clicks on any of the nodes of the tree, the corresponding visual annotations are triggered. Thus in this case if sphere was mapped into a Jawaa “circle” using the annotation: sphere(A,B)==>circle##name=spherevalueOf(arg0)#x=30#y=prologValue(R is random(30), X is R*15)#width=30#color=black#bkgrd=valueOf(arg1). If the user clicks on the node with the constraints sphere(b,blue) and sphere(a,green), the system automatically connects the constraints to the previously introduced visual tracer that checks if any of the current constraints have annotations. This produces a Jawaa animation file with two circles placed randomly. An example of the visualization result is shown in Figure 8.

7 Conclusion

The paper introduced a new tool that is able to visualize different CHR programs by dynamically linking CHR constraints to visual objects. To overcome the problems faced with the previous tools, the new system outsources the visualization process to existing tools to have a generic tracing technique. Intelligence is thus shifted to the transformation and annotation modules. Through the provided set of visual objects and actions, different algorithms could be annotated to produce the needed visual states. Such visualization features have proved to be useful in
many situations including code debugging for programmers and educational purposes [17].

In addition the paper also explores the possibility of visualizing the execution of different operational semantics of \( \text{CHR} \). It provides a module that is able to visualize the exhaustive execution of \( \text{CHR} \) and more importantly it links it to the annotated constraints.

Thus in general, unlike the previously provided tools [18] for visualizing constraint programs, the focus is not just on the search space and the domains. The provided tool enables its users to focus on the algorithms executed to visualize their states.

In the future, more dynamic annotation options could be provided to the user. The visualization of the execution of different \( \text{CHR} \) operational semantics should be investigated. The tool could also be extended to be a visual confluence checker for \( \text{CHR} \) programs. The end goal is to have a generic source-to-source transformation workbench for \( \text{CHR} \).
References

Declarative Compilation for Constraint Logic Programming

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Abstract. We present a new declarative compilation of logic programs with constraints into variable-free relational theories which are then executed by rewriting. This translation provides an algebraic formulation of the abstract syntax of logic programs. Management of logic variables, unification, and renaming apart is completely elided in favor of algebraic manipulation of variable-free relation expressions. We prove the translation is sound, and the rewriting system complete with respect to traditional SLD semantics.

Keywords: logic programming, constraint programming, relation algebra, rewriting, semantics

1 Introduction

Logic programming is a paradigm based on proof search and directly programming with logical theories. This is done to achieve declarative transparency: guaranteeing that execution respects the mathematical meaning of the program.

The power that such a paradigm offers comes at a cost for formal language research and implementation. Management of logic variables, unification, renaming variables apart and proof search are cumbersome to handle formally. Consequently, it is often the case that the formal definition of these aspects is left outside the semantics of programs, complicating reasoning about them and the introduction of new declarative features.

We address this problem here by proposing a new compilation framework — based on ideas of Tarski\textsuperscript{22} and Freyd\textsuperscript{9} — that encodes logic programming syntax into a variable-free algebraic formalism: relation algebra.

Relation algebras are pure equational theories of structures containing the operations of composition, intersection and convolution. An important class of relation algebras is the so-called distributive relation algebras with quasi-projections, which also incorporate union and projections.

We present the translation of constraint logic programs to such algebras in 3 steps. First, for a CLP program $P$ with signature $\Sigma$, we define its associated relation algebra $\text{QRA}_P$, which provides both the target object language for program translation and formal axiomatization of constraints and logic variables.
Second, we introduce a constraint compilation procedure that maps constraints to variable-free relation terms in $\text{QRA}_\Sigma$. Third, a program translation procedure compiles constraint logic programs to an equational theory over $\text{QRA}_\Sigma$.

The key feature of the semantics and translation is its variable-free nature. Programs that contain logical variables are represented as ground terms in our setting, thus all reasoning and execution is reduced to algebraic equality, allowing the use of rewriting. The resulting system is sound and complete with respect to SLD resolution.

Our compilation provides a solution to the following problems:

– Underspecification of abstract syntax and logic variable management in logic programs: solved by the inclusion of metalogical operations directly into the compilation process.
– Interdependence of compilation and execution strategies: solved by making target code completely orthogonal to execution.
– Lack of transparency in compilation (for subsequent optimization and abstract interpretation): solved by making target code a low-level yet fully declarative translation of the original program.

Variable elimination and relation composition

We illustrate the spirit of translation, and in particular the variable elimination procedure, by considering a simple case, namely the transitive closure of a graph:

$$\text{edge}(a,b). \quad \text{connected}(X,X).$$
$$\text{edge}(b,c). \quad \text{connected}(X,Y) :- \text{edge}(X,Z), \text{connected}(Z,Y).$$
$$\text{edge}(a,e).$$
$$\text{edge}(e,f).$$

In this carefully chosen example, the elimination of variables and the translation to binary relation symbols is immediate:

$$\text{edge} = (a,b) \cup (b,c) \cup (a,e) \cup (a,e) \cup (e,f)$$
$$\text{connected} = \text{id} \cup \text{edge} : \text{connected}$$

The key feature of the resulting term is the composition $\text{edge} : \text{connected}$. The logical variable $Z$ is eliminated by the composition of relations allowing the use of variable free object code.

A query $\text{connected}(a,X)$ is then modeled by the relation $\text{connected} \cap (a,a) \text{1}$ where $\text{1}$ is the (maximal) universal relation. Computation can proceed by rewriting the query using a suitable orientation of the relation algebra equations and unfolding pertinent recursive definitions.

Handling actual arbitrary constraint logic programs is more involved. First, we use projection relations to handle predicates involving an arbitrary number of arguments and an unbounded number of logic variables; second, we formalize constraints in a relational way.
Projections and permutations algebraically encode all the operations of logical variables, disjunctive and conjunctive clauses are handled with the help of the standard relational operators $\cap$, $\cup$.

**Constraint Logic Programming Conventions** We refer the reader to [16] for basic definitions of logic programming over Horn clauses, and [12] for background on the syntax and semantics of constraint logic programming. In this paper we fix a signature $\Sigma$, a set of terms $T_\Sigma(X)$, and a subset $C$ of all first-order formulas over $\Sigma$ closed under conjunction and existential quantification to be the set of constraint formulas as well as a $\Sigma$-structure $D$, called the constraint domain. Constraint logic programs are sets of Horn clauses. We use vector notation extensively in the paper, to abbreviate Horn clauses with constraints $p \leftarrow q_1, \ldots, q_n$, where $p$ is an atomic formula and $q_i$ may be an atomic formula or a constraint. For instance, in our vector notation, a clause is written $p(t[x]) \leftarrow q(u[x, y])$, where the boldface symbols indicate vectors of variables $x, y$, terms $t, u$ (depending on variables $x$, etc...) and predicates $q$ (depending on terms $u$).

## 2 Relation Algebras and Signatures

In this section, we define $QRA_\Sigma$, a relation algebra in the style of [22,9] formalizing a CLP signature $\Sigma$ and a constraint domain $D$. We define its language, its equational theory and semantics.

### 2.1 Relational Language and Theory

The relation language $R_\Sigma$ is built from a set $R_C$ of relation constants for constant symbols, a set $R_F$ of relation constants for function symbols from $\Sigma$, and a set of relation constants for primitive predicates $R_{CP}$, as well as a fixed set of relation constants and operators detailed below. Let us begin with $R_C$. Each constant symbol $a \in C_\Sigma$ defines a constant symbol $(a, a) \in R_C$, each function symbol $f \in F_\Sigma$ defines a constant symbol $R_f$ in $R_F$. Each predicate symbol $r \in CP_\Sigma$ defines a constant symbol $r$ in $R_{CP}$. We write $R_\Sigma$ for the full relation language:

$$
R_C = \{(a, a) \mid a \in C_\Sigma\} \quad R_F = \{f \mid f \in F_\Sigma\} \quad R_{CP} = \{r \mid r \in CP_\Sigma\}
$$

$$
R_{atom} ::= R_C \mid R_F \mid R_{CP} \mid id \mid di \mid 1 \mid 0 \mid hd \mid tl
$$

$$
R_{\Sigma} ::= R_{atom} \mid R_{\Sigma}^\circ \mid R_{\Sigma} \cup R_{\Sigma} \mid R_{\Sigma} \cap R_{\Sigma} \mid R_{\Sigma} R_{\Sigma}
$$

The constants $1, 0, id, di$ respectively denote the universal relation (whose standard semantics is the set of all ordered pairs on a certain set), the empty relation, the identity (diagonal) relation, and identity’s complement. Juxtaposition $RR$ represents relation composition (often written $R;R$) and $R^\circ$ is the inverse of $R$. We write $hd$ and $tl$ for the head and tail relations. The projection of an $n$-tuple onto its $i$-th element is written $P_i$ and defined as $P_1 = hd, P_2 = tl; hd, \ldots, P_n = tl^{n-1}; hd.

$QRA_\Sigma$ (Fig. [1]) is the standard theory of distributive relation algebras, plus Tarski’s quasiprojections [22], and equations axiomatizing the new relations of $R_{\Sigma}$.
Note that products and their projections are axiomatized in a relational, variable-free manner.

2.2 Semantics

Let $\Sigma$ be a constraint signature and $D$ a $\Sigma$-structure. Write $t^D$ for the interpretation of a term $t \in T_\Sigma$. We define $D^\dagger$ to be the union of $D^0 = \{\langle \rangle\}$ (the empty sequence), $D$ and $D$-finite products, for example: $D^2, D^2 \times D, D \times D^2, \ldots$. We write $\langle a_1, \ldots, a_n \rangle$ for members of the $n$-fold product associating to the right, that is to say, $\langle a_1, \langle a_2, \ldots, \langle a_{n-1}, a_n \rangle \cdots \rangle \rangle$. Furthermore, we assume right-association of products when parentheses are absent. Note that the 1 element sequence does not exist in the domain, so we write $\langle a \rangle$ for $a$ as a convenience.

Let $R_D = D^\dagger \times D^\dagger$. We make the power set of $R_D$ into a model of the relation calculus by interpreting atomic relation terms in a certain canonical way, and the operators in their standard set-theoretic interpretation. We interpret $hd$ and $tl$ as projections in the model.

**Definition 1.** Given a structure $D$ a relational $D$-interpretation is a mapping $[\_]^D\dagger$ of relational terms into $R_D$ satisfying the identities in Fig. 2. The function $\alpha$ used in this table and elsewhere in this paper refers to the arity of its argument, whether a relation or function symbol from the underlying signature.

**Theorem 1.** Equational reasoning in QRA$_\Sigma$ is sound for any interpretation:

$$\vdash_{QRA} R = S \implies [R]^D\dagger = [S]^D\dagger$$
We fix a canonical list of variables to a variable-free relational term. For instance, the term $x_1 + s(s(x_2))$ is translated to the relation $\bigcup_{i \leq n} P_i \cap P_2; s; s; P_2^2$.

Fig. 2. Standard interpretation of binary relations.

3 Program Translation

We define constraint and program translation to relation terms. To this end, we define a function $K$ from constraint formulas with — possibly free — logic variables to a variable-free relational term. $K$ is the core of the variable elimination mechanism and will appear throughout the rest of the paper.

The reader should keep in mind that there are two kinds of predicate symbols in a constraint logic program: constraint predicates $r$ which are translated by the function $K$ above to relation terms $r$, and defined or program predicates.

We translate defined predicates — and CLP programs — to equations $\bar{p} \equiv R$, where $\bar{p}$ will be drawn from a set of definitional variables standing for program predicate names $p$, and $R$ is a relation term. The set of definitional equations can be both seen as an executable specification, by understanding it in terms of the rewriting rules given in this paper; or as a declarative one, by unfolding the definitions and using the standard set-theoretic interpretation of binary relations.

3.1 Constraint Translation

We fix a canonical list $x_1, \ldots, x_n$ of variables occurring in all terms, so as to translate them to variable-free relations in a systematic way. There is no loss of generality as later, we transform programs into this canonical form.

**Definition 2 (Term Translation).** We define a translation function $K : T_\Sigma(\mathcal{X}) \to \mathcal{R}_\Sigma$ from first-order terms to relation expressions as follows:

\[
K(c) = (c, c) \mathbf{1} \\
K(x_i) = P^P \\
K(f(t_1, \ldots, t_n)) = R_f; \bigcap_{i \leq n} P_i; K(t_i)
\]

This translation is extended to vectors of terms as follows $K((t_1, \ldots, t_n)) = \bigcap_{i \leq n} P_i; K(t_i)$.

The semantics of the relational translation of a term is the set of all of the instances of that term, paired with the corresponding instances of its variables. For instance, the term $x_1 + s(s(x_2))$ is translated to the relation $+; (P_1; P_1^P \cap P_2; s; s; P_2^2)$. 
Lemma 1. Let $t[x]$ be a term of $T_\Sigma(\mathcal{X})$ whose free variables are among those in
the sequence $x = x_1, \ldots, x_m$. Then, for any sequences $a = a_1, \ldots, a_m \in D^\uparrow, u \in D^\uparrow$ and any $b \in D$ we have

$$(b, au) \in [K(t[x])]^D \iff b \in t^D[a/x]$$

We will translate constraints over $m$ variables to partially coreflexive relations
over the elements that satisfy them. A binary relation $R$ is coreflexive if it is
contained in the identity relation, and it is $i$-coreflexive if its $i$-th projection is
contained in the identity relation: $P_i^\circ; R; P_i \subseteq id$. Thus, for a variable $x_i$ free in
a constraint, the translation will be $i$-coreflexive.

We now formally define two partial identity relation expressions $I_m, Q_i$ for the
translation of existentially quantified formulas, in such a way that if a constraint
$\varphi[x]$ over $m$ variables is translated to an $m$-coreflexive relation, the formula
$\exists x_i. \varphi[x]$ corresponds to a coreflexive relation in all the positions but the $i$-th
one, as $x_i$ is no longer free. In this sense $Q_i$ may be seen as a hiding relation.

Definition 3. The partial identity relation expressions $I_m, Q_i$ for $m, i > 0$ are
defined as:

$$I_m := \bigcap_{1 \leq i \leq m} P_i(P_i)^\circ \quad Q_i = I_{i-1} \cap J_{i+1} \quad J_i = tl^i; (tl^i)^i$$

$I_m$ is the identity on sequences up to the first $m$ elements. $Q_i$ is the identity on
all but the $i$-th element, with the $i$-th position relating arbitrary pairs of elements.

Definition 4 (Constraint Translation). The $\hat{K} : \mathcal{L}_D \to \mathcal{R}_\Sigma$ translation func-
tion for constraint formulas is:

$$\begin{align*}
\hat{K}(p(t_1, \ldots, t_n)) &= \bigcap_{1 \leq i \leq n} K(t_i)^\circ; P_i; \bigcap_{1 \leq i \leq n} P_i; K(t_i) \\
\hat{K}(\varphi \land \theta) &= \hat{K}(\varphi) \cap \hat{K}(\theta) \\
\hat{K}(\exists x_i. \varphi) &= Q_i; \hat{K}(\varphi); Q_i
\end{align*}$$

As an example, the translation of the constraint $\exists x_1, x_2. s(x_1) \leq x_2$ is

$$Q_1; Q_2; (P_1^\circ; s^\circ; P_1 \cap P_2^\circ; P_2); \leq; (P_1; s; P_1^\circ \cap P_2; P_2^\circ); Q_1; Q_2$$

Lemma 2. Let $\varphi[x]$ be a constraint formula with free variables among $x =
\begin{align*}
&\begin{align*}
(x_1, \ldots, x_m). \quad &\text{Then, for any sequences } a = a_1, \ldots, a_m, u &\text{ and } u' \text{ of members of } D \\
&(a u, a u') \in \hat{K}(\varphi[x]) \iff D \models \varphi[a/x]
\end{align*}
\end{align*}$$

3.2 Translation of Constraint Logic Programs

To motivate the technical definitions below, we illustrate the program translation
procedure with an example. Assume a language with constant 0, a unary function
symbol $s$, constraint predicate $=$ and program predicate $add$. We can write the
traditional Horn clause definition of Peano addition:
The clause’s GPF form is:

\[ \text{add}(0, X, X) \cdot \text{add}(s(X), Y, s(Z)) :\text{add}(X, Y, Z) \cdot \]

This program is first purified: the variables in the head of the clauses defining each predicate are chosen to be a sequence of fresh variables \(x_1, x_2, x_3\), with all bindings stated as equations in the tail.

\[ \text{add}(x_1, x_2, x_3) \leftarrow x_1 = 0, x_2 = x_3. \]
\[ \text{add}(x_1, x_2, x_3) \leftarrow \exists x_4, x_5 x_1 = s(x_4), x_3 = s(x_5), \text{add}(x_4, x_2, x_5) \]

The clauses are combined into a single definition similar to the Clark completion of a program. We also use the variable permutation \(\pi\) sending \(x_1, x_2, x_3, x_4, x_5 \mapsto x_4, x_2, x_5, x_1, x_3\) to rewrite the occurrence of the predicate \text{add}\ in the tail so that its arguments coincide with those in the head.

\[ \text{add}(x_1, x_2, x_3) \leftrightarrow (x_1 = 0, x_2 = x_3) \]
\[ \lor \exists x_4, x_5 x_1 = s(x_4), x_3 = s(x_5), \pi \text{add}(x_1, x_2, x_3). \]

Now we apply relational translation \(\hat{K}\) defined above to all relation equations, and eliminate the existential quantifier using the partial identity operator \(I_3\) defined below. We represent the permutation \(\pi\) using the relation expression \(W_\pi\) that simulates its behavior in a variable-free manner as described below, and replace the predicate \text{add}\ with a corresponding relation variable \(\overline{\text{add}}\).

\[ \overline{\text{add}} = \hat{K}(x_1 = o \land x_2 = x_3) \cup I_3((\hat{K}(x_1 = s(x_4) \land x_3 = s(x_5)) \cap W_\pi(\overline{\text{add}}))) \]

Now we give a description of the general translation procedure. We first process programs to their complete database form as defined in [6], which given the executable nature of our semantics reflects the choice to work within the minimal semantics. The main difference in our processing of a program \(P\) to its completed form \(P'\) is that a strict policy on variable naming is enforced, so that the resulting completed form is suitable for translation to relational terms.

**Definition 5 (General Purified Form for Clauses).** For a clause \(p(t[y]) \leftarrow q(v[y]),\) let \(h = \alpha(p), y = [y], v = [v],\) and \(m = h + y + v.\) Assume vectors:

\[
x = x_h x_t x_y x_v = x_1, \ldots, x_h, x_{h+1}, \ldots, x_{h+y}, x_{h+y+1}, \ldots, x_m \\
x_h = x_1, \ldots, x_h \\
x_t = x_y x_v = x_{h+1}, \ldots, x_{h+y}, x_{h+y+1}, \ldots, x_m \\
x_y = x_{h+1}, \ldots, x_{h+y} \\
x_v = x_{h+y+1}, \ldots, x_m
\]

the clause’s GPF form is:

\[ p(x_h) \leftarrow \exists^{h\uparrow} \left((x_h = t[x_y] \land x_v = v[x_y]), q(x_v) \right) \]

\(\exists^n\) denotes existential closure with respect to all variables whose index is greater than \(n.\) \(x_h\) and \(x_t\) stand for head and tail variables. A program is in GPF form iff every one of its clauses is. After the GPF step, we perform Clark’s completion.
Definition 6 (Completion of a Predicate). We define Clark’s completed form for a predicate $p$ with clauses $cl_1, \ldots, cl_n$ in GPF form:

$$
p(x_h) \leftarrow cl_1, tl_1
\ldots
p(x_h) \leftarrow cl_n, tl_k
\Rightarrow
p(x_h) \leftrightarrow tl_1 \lor \cdots \lor tl_k
$$

The above definition easily extends to programs. Completed forms are translated to relations by using $\hat{K}$ for the constraints, mapping conjunction to $\cap$ and $\lor$ to $\cup$. Existential quantification, recursive definitions and parameter passing are handled in a special way which we proceed to detail next.

Existential Quantification: Binding Local Variables Variables local to the tail of a clause are existentially quantified. For technical reasons — simpler rewrite rules — we use the partial identity relation $I_n$, rather than the $Q_n$ relation defined in the previous sections. $I_n$ acts as an existential quantifier for all variables of index greater than a given number.

Lemma 3. Let $a = a_1, \ldots, a_n \in D$, $x = x_1, \ldots, x_n$, let $\varphi$ be a constraint over $m$ free variables, with $m > n$, $y$ a vector of length $k$ such that $n + k = m$, and $u, v \in D^*$, then:

$$(au, av) \in [I_n; \hat{K}(\varphi[xy]); I_n]^D \iff D \models (\exists^{n+k} \varphi[xy])[a/x]$$

Recursive Predicate Definitions We shall handle recursive predicate definitions by extending the relational language with a set of definitional symbols $\bar{p}, \bar{q}, \bar{r}, \ldots$ for predicates. Then, a recursive predicate $\bar{p}$ is translated to a definitional equation $\bar{p} \equiv R(\bar{p}_1, \ldots, \bar{p}_n)$, spelled out in Def. 8 where the notation $R(\bar{p}_1, \ldots, \bar{p}_n)$ indicates that relation $R$ resulting from the translation may depend on predicate symbols $\bar{p}_1, \ldots, \bar{p}_n$. Note that $R$ is monotone in $\bar{p}_1, \ldots, \bar{p}_n$. Consequently, using a straightforward fixed point construction we can extend the interpretation $[\_]^D$ to satisfy $[\bar{p}]^D = [R(\bar{p}_1, \ldots, \bar{p}_n)]^D$, thus preserving soundness when we adjoin the definitional equations to $\text{QRA}_\Sigma$.

Parameter Passing The information about the order of parameters in each pure atomic formula $p(x_{i_1}, \ldots, x_{i_r})$ is captured using permutations. Given a permutation $\pi : \{1..n\} \rightarrow \{1..n\}$, the function $w_\pi$ on formulas and terms is defined in the standard way by its action over variables. We write $W_\pi$ for the corresponding relation:

Definition 7 (Switching Relations). Let $\pi : \{1..n\} \rightarrow \{1..n\}$ be a permutation. The switching relation expression $W_\pi$, associated to $\pi$ is:

$$W_\pi = \bigcap_{j=1}^{n} P_{\pi(j)}(P_j)^{\circ}.$$ 

Lemma 4. Fix a permutation $\pi$ and its corresponding $w_\pi$ and $W_\pi$. Then:

$$[\hat{K}(w_\pi(p(x_1, \ldots, x_n)))] = [W_\pi \hat{K}(p)W_\pi^\circ]$$
Declarative Compilation for Constraint Logic Programming

mother(sarah, isaac).
parent(X, Y) ← father(X, Y).
parent(X, Y) ← mother(X, Y).
sibling(S1, S2) ← S1 ≠ S2, parent(Par, S1), parent(Par, S2).
brother(Brother, Sib) ← male(Brother), sibling(Brother, Sib).

Fig. 3. Biblical family relations in Prolog.

The Translation Function

Now we may define the translation for defined predicates.

Definition 8 (Relational Translation of Predicates). Let $h, p(x_h)$ be as in Def. 5. The translation function $Tr$ from completed predicates to relational equations is defined by:

\[\begin{align*}
Tr(p(x_h) \leftarrow c_1 \lor \cdots \lor c_k) &= (p \models \bigcup Tr(c_1) \cup \cdots \cup Tr(c_k)) \\
Tr_{cl}(\exists x^t, p) &= I_h; (Tr_l(p_1) \cap \cdots \cap Tr_l(p_n)); I_h \\
Tr_l(\varphi) &= \mathcal{K}(\varphi) \quad \varphi \text{ a constraint} \\
Tr_l(p_i(x_i)) &= W^{\pi}; \mathcal{P}_i; W^\varnothing_i \text{ such that } \pi(x_i) = x_1, \ldots, x_{\alpha(p_i)}
\end{align*}\]

Example 1. Figure 3 shows a fragment of a constraint logic program to represent a family relations database [21]. Consider the translation of the program predicates mother, parent, sibling and brother. We write the program in general purified form:

\[
\begin{align*}
\text{mother}(x_1, x_2) &\iff (x_1 = \text{sarah}) \land (x_2 = \text{isaac}) \\
\text{parent}(x_1, x_2) &\iff \text{father}(x_1, x_2) \lor \text{mother}(x_1, x_2) \\
\text{sibling}(x_1, x_2) &\iff \exists x_3. x_1 \neq x_2 \land \text{parent}(x_3, x_1) \land \text{parent}(x_3, x_2) \\
\text{brother}(x_1, x_2) &\iff \text{male}(x_1) \land \text{sibling}(x_1, x_2)
\end{align*}
\]

Letting $\sigma_1$ and $\sigma_2$ be the permutations $\langle 1, 2, 3 \rangle \rightarrow \langle 3, 1, 2 \rangle$ and $\langle 1, 2, 3 \rangle \rightarrow \langle 3, 2, 1 \rangle$ respectively we obtain

\[
\begin{align*}
\text{mother} &= \mathcal{K}(x_1 = \text{sarah}) \cap \mathcal{K}(x_2 = \text{isaac}) \\
\text{parent} &= \text{father} \cup \text{mother} \\
\text{sibling} &= \mathcal{K}(x_1 \neq x_2) \cap \mathcal{I}_2[W_{\sigma_1}\text{parent}W^\varnothing_{\sigma_1} \cap W_{\sigma_2}\text{parent}W^\varnothing_{\sigma_2}]; \mathcal{I}_2 \\
\text{brother} &= \text{male} \cap \text{sibling}
\end{align*}
\]
The query `brother(X,milcah)` leads to the rewriting of the term `\( \dot{K}(x_2 = \text{milcah}) \cap \text{brother} \)` to `\( \dot{K}(x_2 = \text{milcah}) \cap \dot{K}(x_1 = \text{lot}) \)`.

4 A Rewriting System for Resolution

In this section, we develop a rewriting system for proof search based on the equational theory \( \text{QRA}_\Sigma \), which will be proven equivalent to the traditional operational semantics for CLP. In Sec. 5 we will show that answers obtained by resolution correspond to answers yielded by our rewriting system and conversely.

The use of ground terms permits the use of rewriting, overcoming the practical and theoretical difficulties that the existence of logic variables causes in equational reasoning. Additionally, we may speak of executable semantics: we use the same function to compile and interpret CLP programs in the relational denotation.

For practical reasons, we don’t rewrite over the full relational language, but we will use a more compact representation of the relations resulting from the translation.

Formally, the signature of our rewriting system is given by the following term-forming operations over the sort \( T_R \): \( I : (\mathbb{N} \times T_R) \rightarrow T_R \), \( W : (\text{Perm} \times T_R) \rightarrow T_R \), \( K : L_D \rightarrow T_R \), \( \cup : (T_R \times T_R) \rightarrow T_R \) and \( \cap : (T_R \times T_R) \rightarrow T_R \). Thus, for instance, the relation \( I_n : R \); \( I_n \) is formally represented in the rewriting system as \( I(n,R) \), provided \( R \) can be represented in it. In practice we make use of the conventional relational notation \( I_n, W_\pi \) when no confusion can arise.

4.1 Meta-reductions

We formalize the interface between the rewrite system and the constraint solver as meta-reductions (Fig. 4). Every meta-reduction uses the constraint solver in a black-box manner to perform constraint manipulation and satisfiability checking.

**Lemma 5.** All meta-reductions are sound: if \( m_i : l \stackrel{r}{\rightarrow} r \) then \( \llbracket l \rrbracket^{D^1} = \llbracket r \rrbracket^{D^1} \).

---

4 There is no problem in defining the rewriting system using the general relational signature, but we would need considerably more rules for no gain.
We prove that our rewriting system over relational terms simulates "traditional" SLD proof search. We present a rewriting system for proof search in Fig. 5. We prove local confluence.

4.2 A Rewriting system for SLD resolution

We present a rewriting system for proof search in Fig. 5. We prove local confluence.

Lemma 6. \( \overset{p}{\rightarrow} \) is sound: if \( p_1 : l \overset{p}{\rightarrow} r \) then \( \llbracket l \rrbracket^D = \llbracket r \rrbracket^D \).

Lemma 7. If we give higher priority to \( p_7 \) over \( p_8 \), \( \overset{p_7}{\rightarrow} \) is locally confluent.

A left outermost strategy gives priority to \( p_7 \) over \( p_8 \).

5 Operational Equivalence

We prove that our rewriting system over relational terms simulates "traditional" SLD proof search specified as a transition-based operational semantics (i.e. [12,7]). For reasons of space, we give a high-level overview of the proof. The full details can be found in the online technical report.

Recall a resolvent is a sequence of atoms or constraints \( p \). We write \( \square \) for the empty resolvent. We assume given a constraint domain \( D \) and its satisfaction relation \( D \models \varphi \). A program state is an ordered pair \( \langle p, \varphi \rangle \) where \( p \) is a resolvent and \( \varphi \) is a constraint (called the constraint store). The notation \( cl : p(u[y]) \leftarrow q(v[z]) \) indicates that \( p(u[y]) \leftarrow q(v[z]) \) is a program clause with label \( cl \). Then, the standard operational semantics for SLD resolution can be defined as the following transition system over program states:

Definition 9 (Standard SLD semantics).

\[
\begin{align*}
\langle \varphi, p \mid \psi \rangle & \xrightarrow{c_{x \leftarrow l}} \langle \psi \mid \varphi \rangle & \text{iff } D \models \varphi \\
\langle p(t[x]), p \mid \varphi \rangle & \xrightarrow{res_{x\leftarrow l}} \langle q(v[\sigma(z)]), p \mid \varphi \wedge (u[\sigma(y)] = t[x]) \rangle \\
\text{where: } cl : p(u[y]) \leftarrow q(v[z]) & \\
D \models \varphi \wedge (u[\sigma(y)] = t[x]) & \\
\sigma \text{ a renaming apart for } y, z, x
\end{align*}
\]
Taking the previous system as a reference, the proof proceeds in two steps: we first define a new transition system that internalizes renaming apart and proof search, and we prove it equivalent to the standard one.

Second, we show a simulation relation between the fully internalized transition system and a transition system defined over relations, which is implemented by the rewriting system of Sec. 4.

With these two equivalences in place, the main theorem is:

**Theorem 2.** The rewriting system of Fig. 3 implements the transition system of Def. 9. Formally, for every transition \((r_1, r_2) \in (\to_l)^*\),

\[ \exists n. (\text{Tr}(r_1), \text{Tr}(r_2)) \in (\mapsto_P)^n \]

and

\[ \forall r_3. (\text{Tr}(r_1), r_3) \in (\mapsto_P)^n \Rightarrow \text{Tr}(r_2) = r_3 \]

Thus, given a program \(P\), relational rewriting of translation will return an answer constraint \(K(\varphi)\) iff SLD resolution from \(P\) reaches a program state \(\langle \square | \varphi' \rangle\), with \(\varphi \iff \varphi'\).

In the next section, we briefly describe the main intermediate system used in the proof.

### 5.1 The resolution transition system

The crucial part of the SLD-simulation proof is the definition of a new extended transition system over program states that will internalize both renaming apart and the proof-search tree. It is an intermediate system between relation rewriting and traditional proof search.

The first step towards the new system is the definition of an extended notion of state. In the standard system of Def. 9, a state is a resolvent plus a constraint store. Our extended notion of state includes:

- A notion of *scope*, which is captured by a natural number which can be understood as the number of global variables of the state.
- A notion of *substate*, which includes information about parameter passing in the form of a *permutation*.
- A notion of clause *selection*, and
- A notion of *failure* and *parallel state*, which represents failures in the search tree and alternatives.

Such states are enough to capture all the meta-theory of constraint logic programming except recursion, which operates meta-logically by replacing predicate symbols by their definitions. Formally:

**Definition 10.** The set \(\mathcal{PS}\) of resolution states is inductively defined as:

- \(\{\text{fail}\}\).
- \(\langle p \mid \varphi \rangle_n\), where \(p_i \equiv P_i(x_i)\) is an atom, or a constraint \(p_i \equiv \psi, x_i\) a vector of variables, \(\varphi\) a constraint store and \(n\) a natural number.
– \( \langle \pi \; PS, \; p \; | \varphi \rangle \rangle_n \), where \( PS \) is a resolution state, and \( \pi \) a permutation.
– \( \langle \pi \; PS, \; p \; | \varphi \rangle \rangle_n \), the “select state”. It represents the state just before selecting a clause to proceed with proof search.
– \( (PS_1 \parallel PS_2) \). The bar is parallel composition, capturing choice in the proof search tree.

The resolution transition system \( \rightarrow_P \subseteq (PS \times PS) \) is shown in Fig. 6. The two first transitions deal with the case where a constraint is first in the resolvent, failing or adding it to the constraint store in case it is satisfiable.

When the head of the resolvent is a defined predicate, the call transition will replace it by its definition, properly encapsulated by a select state equipped with the permutation capturing argument order.

The select transition performs two tasks: first, it modifies the current constraint store adding the appropriate permutation and scoping \( (n, \pi) \); second, it selects the first clause for proof search.

The return transitions will either propagate failure or undo the permutation and scoping performed at call time.

\( sub, \) backtrack, and seq are structural transitions with a straightforward interpretation from a proof search perspective.

Then, we have the following lemma:

Fig. 6. Resolution Transition System
Lemma 8. For all queries $⟨p | ϕ⟩_n$, the first successful $→_l$ derivation using an SLD strategy uniquely corresponds to a $→_p$ derivation:

$$⟨p | ϕ⟩_n →_l ... →_l ⟨□ | ϕ'⟩_n \iff ⟨p | ϕ⟩_n →_p ... →_p (⟨□ | ϕ'⟩_n \mid PS)$$

for some resolution state $PS$.

Corollary 1. The transition systems of Def. 9 and Fig. 6 are answer-equivalent: for any query they return the same answer constraint.

With this lemma in place, the proof of Thm. 2 is completed by showing a simulation between the resolution system and a transition system induced by relation rewriting.

6 Related and Future Work

Previous work: The paper is the continuation of previous work in [4,15,11] considerably extended to include constraint logic programming, which requires a different translation procedure and a different rewriting system.

In particular, the presence of constraints in this paper permits a different translation of the Clark completion of a program and plays a crucial role in the proof of completeness, which was missing in earlier work. The operational semantics is also new.

Related work: A number of solutions have been proposed to the syntactic specification problem. There is an extensive literature treating abstract syntax of logic programming (and other programming paradigms) using encodings in higher-order logic and the lambda calculus [19], which has been very successful in formalizing the treatment of substitution, unification and renaming of variables, although it provides no special framework for the management and progressive instantiation of logic variables, and no treatment of constraints. Our approach is essentially orthogonal to this, since it relies on the complete elimination of variables, substitution, renaming and, in particular, existentially quantified variables. Our reduction of management of logic variables to variable free rewriting is new, and provides a complete solution to their formal treatment.

An interesting approach to syntax specification is the use of nominal logic [23] in logic programming, another, the formalization of logic programming in categorical logic [2,20,13,1,8] which provides a mathematical framework for the treatment of variables, as well as for derivations [13]. None of the cited work gives a solution that simultaneously includes logic variables, constraints and proof search strategies however.

Bellia and Occhiuto [3] have defined a new calculus, the C-expression calculus, to eliminate variables in logic programming. We believe our translation into the well-understood and scalable formalism of relations is more applicable to extensions of logic programming. Furthermore the authors do not consider constraints.
Future work: A complementary approach to this work is the use of category theory, in particular the Freyd’s theory of tabular allegories [9] which extends the relation calculus to an abstract category of relations providing native facilities for generation of fresh variables and a categorical treatment of monads. A first attempt in this direction has been published by the authors in [10]. It would be interesting to extend the translation to hereditarily Harrop or higher order logic [18] by using a stronger relational formalism, such as Division and Power Allegories. Also, the framework would yield important benefits if it was extended to include relation and set constraints explicitly.

7 Conclusion

We have developed a declarative relational framework for the compilation of Constraint Logic programming that eliminates logic variables and gives an algebraic treatment of program syntax. We have proved operational equivalence to the classical approach. Our framework has several significant advantages.

Programs can be analyzed, transformed and optimized entirely within this framework. Execution is carried out by rewriting over relational terms. In these two ways, specification and implementation are brought much closer together than in the traditional logic programming formalism.

References


Pre-Indexed Terms for Prolog*

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Abstract. Indexing of terms and clauses is a well-known technique used in Prolog implementations (as well as automated theorem provers) to speed up search. In this paper we show how the same mechanism can be used to implement efficient reversible mappings between different term representations, which we call \textit{pre-indexings}. Based on user-provided term descriptions, these mappings allow us to use more efficient data encodings internally, such as prefix trees. We show that for some classes of programs, we can drastically improve the efficiency by applying such mappings at selected program points.

1 Introduction

Terms are the most important data type for languages and systems based on first-order logic, such as (constraint) logic programming or resolution-based automated theorem provers. Terms are inductively defined as variables, atoms, numbers, and compound terms (or structures) comprised by a functor and a sequence of terms.\textsuperscript{3} Two main representations for Prolog terms have been proposed. Early Prolog systems, such as the Marseille and DEC-10 implementations, used \textit{structure sharing} \cite{2}, while the WAM \cite{15,1} –and consequently most modern Prolog implementations– uses \textit{structure copying}. In structure sharing, terms are represented as a pair of pointers, one for the structure skeleton, which is shared among several instances, and another for the binding environment, which determines a particular instantiation. In contrast, structure copying makes a copy of the structure for each newly created term. The encoding of terms in memory resembles tree-like data structures.

In order to speed up resolution, sophisticated term indexing has been implemented both in Prolog \cite{1,7} and automated theorem provers \cite{6}. By using specialized data structures (such as, e.g., tries), indexing achieves sub-linear complexity in clause selection. Similar techniques are used to efficiently store predicate solutions in tabling \cite{13}. This efficient indexing machinery is often

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\textsuperscript{3} Additionally, many Prolog systems implement an extension mechanism for variable domains using \textit{attributed variables}. 
attractive for storing and manipulating program data, such as dynamic predicates. Indexed dynamic predicates offer the benefits of efficient key-value data structures while hiding the implementation details from the user program.

Modulo some issues like variable sharing, there is thus a duality in programming style between explicitly encoding data as terms or encoding data implicitly as tuples in dynamic predicates. However, although both alternatives have some declarative flavor, it is also frequent to find code where, for performance reasons, the data is represented in the end in a quite unnatural way. E.g., the set \(\{1, 2, 3, \ldots, n\}\) can be represented naturally as the term \([1,2,3,\ldots,n]\) (equivalent to a linked list). However, depending on the lifetime and operations to be performed on the data, binary trees, some other map-like structure, or dynamic predicates may be preferable. These changes in representation often propagate through the whole program.

The goal of this paper is to study the merits of term indexing during term creation rather than at clause selection time. We exploit the fact that data has frequently a fixed skeleton structure, and introduce a mapping in order to index and share that part. This mapping is derived from program declarations specifying term encoding (called rtypes, for representation types) and annotations defining the program points where pre-indexing of terms is performed. This is done on top of structure copying, so that no large changes are required in a typical Prolog runtime system. Moreover, the approach does not require large changes in program structure, which makes rtypes easily interchangeable.

We have implemented a prototype as a Ciao package that deals with rtype declarations as well as with some additional syntactic sugar that we provide for marking pre-indexing points.

2 Background

We follow the definitions and naming conventions for term indexing of [3,6]. Given a set of terms \(\mathcal{L}\) (the indexed terms), a binary relation \(R\) over terms (the retrieval condition), and a term \(t\) (the query term), we want to identify the subset \(\mathcal{M} \subseteq \mathcal{L}\) consisting of all the terms \(l\) such that \(R(l, t)\) holds (i.e., such that \(l\) is \(R\)-compatible with \(t\)). We are interested in the following retrieval conditions \(R\) (where \(\sigma\) is a substitution):

- \(\text{unif}(l, t) \iff \exists \sigma \ l\sigma = t\sigma\) (unification)
- \(\text{inst}(l, t) \iff \exists \sigma \ l = t\sigma\) (instance check)
- \(\text{gen}(l, t) \iff \exists \sigma \ l\sigma = t\) (generalization check)
- \(\text{variant}(l, t) \iff \exists \sigma \ l\sigma = t\) and \(\sigma\) is a renaming substitution (variant check)

Example 1. Given \(L = \{h(f(A)), h(f(B, C)), h(g(D))\}\), \(t = h(f(1))\), and \(R = \text{unif}\), then \(M = \{h(f(A))\}\).
The objective of term indexing is to implement fast retrieval of candidate terms. This is done by processing the indexed set \( L \) into specialized data structures (index construction) and modifying this index when terms are inserted or deleted from \( L \) (index maintenance).

When the retrieval condition makes use of the function symbols in the query and indexed terms, it is called function symbol based indexing.

In Prolog, indexing finds the set of program clauses such that their heads unify with a given literal in the goal. In tabled logic programming, this is also interesting for detecting if a new goal is a variant or subsumed by a previously evaluated subgoal [5,12].

Limitations of indexing. Depending on the part of the terms that is indexed and the supporting data structure, the worst case cost of indexing is proportional to the size of the term. When computing hash keys, the whole term needs to be traversed (e.g., computing the key for \( h(f(A)) \) requires walking over \( h \) and \( f \)). This may be prohibitively costly, not only in the maintenance of the indices, but also in the lookup. As a compromise many systems rely only on first argument, first level indexing (with constant hash table lookup, relying on linear search for the selected clauses). However, when the application needs stronger, multi-level indexing, lookup costs are repeated many times for each clause selection operation.

### 3 Pre-indexing

The goal of pre-indexing is to move lookup costs to term building time. The idea that we propose herein is to use a bijective mapping between the standard and the pre-indexed representations of terms, at selected program points. The fact that terms can be partially instantiated brings in a practical problem, since binding a variable may affect many precomputed indices (e.g., precomputed indices for \( H=h(X), G=g(X) \) may need a change after \( X=1 \)). Our proposed solution to this problem is to restrict the mapping to terms of a specific form, based on (herein, user-provided) instantiation types.

**Definition 1 (Instantiation type).** We say that \( t \) is an instance of an instantiation type \( \tau \) (defined as a unary predicate), written as \( \text{check}(\tau(t)) \), if there exists a term \( l \) such that \( \tau(l) \) is in the model of \( \tau \) and \( \text{gen}(l,t) \) (or \( \text{inst}(t,l) \)).

For conciseness, we will describe the restricted form of instantiation types used herein using a specialized syntax “:- \text{rtype Name} \rightarrow \text{Cons}_1 ; \ldots ; \text{Cons}_n”, where each \( \text{Cons}_i \) is a term constructor. Each argument of a term constructor may be another rtype name, or the special values \text{any} (that represents any term or variable) and \text{nv} (that represents any nonvar term). E.g.,:

\[ \ldots \]

\[ 4 \] Despite the syntax being similar to that described in [10], note that the semantics is not equivalent.
The rule above thus corresponds to the predicate:

\[
\text{lst([], []).} \\
\text{lst([_|Xs]) :- lst(Xs).}
\]

**Example 2.** According to the definition above for \( \text{lst} \), the terms \([1,2,3]\) and \([0,2]\) belong to \( \text{lst} \) while \([1|0]\) does not. If \( \text{nv} \) were used instead of \( \text{any} \) in the definition above then \([0,2]\) would also not belong to \( \text{lst} \).

**Type-based pre-indexing.** The idea behind pre-indexing is to maintain specialized indexing structures for each \( \text{rtype} \) (which in this work is done based on user annotations). Conceptually, the indexing structure will keep track of the \( \text{rtype} \) inhabitants constructed during the execution dynamically, assigning a unique identifier (the pre-index key) to each of them. E.g., for \( \text{lst} \) we could assign \( \{[], k_0, [_], k_1, [_, _], k_2, \ldots\} \) (that is, \( k_i \) for each list of length \( i \)).

Translation between pre-indexed and non-pre-indexed forms is defined in terms of a pre-indexing casting. Given a term \( t \) so that \( \text{check}(\tau(t)) \), then \( \exists l \in |\tau| \) (set of “weakest” terms for which \( \tau \) holds) such that \( \text{gen}(l, t) \).

**Definition 2 (Pre-indexing cast).** A pre-indexing cast of type \( \tau \) is a bijective mapping between subsets of terms, denoted by \( \#\tau \), with the following properties:

- For every term \( x \) so that \( \text{check}(\tau(x)) \) (which defines the domain of the mapping), and substitution \( \sigma \), then \( \#\tau(x\sigma) = \#\tau(x)\sigma \) (\( \sigma \)-commutative), and
- the first-level functor of \( \#\tau(x) \) encodes the (indexed) structure of the arguments (so that it uniquely identifies the \( \text{rtype} \) inhabitant).

Informally, the first property ensures that pre-indexing casts can be selectively introduced in a program (whose terms are instantiated enough) without altering the (substitution) semantics. Moreover, the meaning of many built-ins is also preserved after pre-indexing, as expressed in the following theorem.

**Theorem 1 (Built-in homomorphism).** Given \( \text{check}(\tau(x)) \) and \( \text{check}(\tau(y)) \), then \( \text{unif}(x, y) \Leftrightarrow \text{unif}(\#\tau(x), \#\tau(y)) \) (equivalently for \( \text{gen}, \text{inst}, \text{variant} \), and other built-ins like \( =\!/2, \text{ground}/1 \)).

**Proof.** \( \text{unif}(x, y) \Leftrightarrow \text{[def. of unif]} \exists \sigma x\sigma = y\sigma \). Since \( \#\tau \) is bijective, then \( \#\tau(x\sigma) = \#\tau(y\sigma) \Leftrightarrow \text{[\( \sigma \)-commutative]} \#\tau(x)\sigma = \#\tau(y)\sigma \). Given the def. of \( \text{unif} \), it follows that \( \text{unif}(\#\tau(x), \#\tau(y)) \). The proofs for other built-ins are similar.

In this work we do not require the semantics of built-ins like \( \&\times \) (i.e., term ordering) to be preserved, but if desired this can be achieved by selecting carefully the order of keys in the pre-indexed term. Similarly, functor arity in principle will not be preserved since ground arguments that are part of the \( \text{rtype} \) structure are allowed to be removed.
3.1 Building pre-indexed terms

We are interested in building terms directly into their pre-indexed form. To achieve this we take inspiration from WAM compilation. Complex terms in variable-term unifications are decomposed into simple variable-structure unifications \( X = f(A_1, \ldots, A_n) \) where all the \( A_i \) are variables. In WAM bytecode, this is further decomposed into a \texttt{put\_str} \( f/n \) (or \texttt{get\_str} \( f/n \)) instruction followed by a sequence of \texttt{unify\_arg} \( A_i \). These instructions can be expressed as follows:

\[
\begin{align*}
\text{put\_str}(X,F/N,S0,S1), & \quad \% | F/N | \\
\text{unify\_arg}(A1,S1,S2) & \quad \% | F/N | A1 | \\
\ldots
\text{unify\_arg}(An,Sn,S) & \quad \% | F/N | A1 | \ldots | An |
\end{align*}
\]

where the \( S_i \) represent each intermediate heap state, which is illustrated in the comments on the right.

Assume that each argument \( A_i \) can be decomposed into its indexed part \( A_i^k \) and its value part \( A_i^v \) (which may omit information present in the key). Pre-indexing builds terms that encode \( A_i^k \) into the main functor by incremental updates:

\[
\begin{align*}
g\text{\_put\_str}(X,F/N,S0,S1), & \quad \% | F/N | \\
g\text{\_unify\_arg}(A1,S1,S2) & \quad \% | F/N<A1^k> | A1^v | \\
\ldots
\text{unify\_arg}(An,Sn,S) & \quad \% | F/N<A1^k,\ldots,An^k> | A1^v | \ldots | An^v |
\end{align*}
\]

The \texttt{rtype} constructor annotations (that we will see in Section 3.2) indicate how the functor and arguments are indexed.

**Cost analysis.** Building and unifying pre-indexed terms have impact both on performance and memory usage. First, regarding time, although pre-indexing operations can be slower, clause selection becomes faster, as it avoids repetitive lookups on the fixed structure of terms. In the best case, \( O(n) \) lookups (where \( n \) is the size of the term) become \( O(1) \). Other operations like unification are sped-up (e.g., earlier failure if keys are different). Second, pre-indexing has an impact on memory usage. Exploiting the data structure allows more compact representations, e.g., \texttt{bitpair(bool,bool)} can be assigned an integer as key (without storage costs). In other cases, the supporting \textit{index structures} may effectively share the common part of terms (at the cost of maintaining those structures).

3.2 Pre-indexing Methods

Pre-indexing is enabled in an \texttt{rtype} by annotating each constructor with modifiers that specify the \textit{indexing method}. Currently we support compact trie-like representations and packaged integer encodings.
Trie representation is specified with the `index(Args)` modifier, which indicates the order in which arguments are walked in the decision-tree. The process is similar to term creation in the heap, but instead of moving a heap pointer, we combine it with walking through a trie of nodes. Keys are retrieved from the term part that corresponds to the `rtype` structure.

Figure 1. Example terms for pre-indexing

For example, let us consider the input set of terms `[a(x), c(z)], [a(x), d(w)], [b(y), c(z)], [b(y), d(w)]`, where `a, b, c, d` are function symbols and `x, y, z, w` are variable symbols. The heap representation is shown in Fig. 1.\footnote{Remember that `[1,2] = .(1,(2,[]))`.} We will compare different `rtype` definitions for representing these terms.

As mentioned before, `nv` represents the `rtype` for any nonvar term (where its first level functor is part of the type). The declaration:

```prolog
:- rtype lst ---> []; [nv|lst]:.:index([0,1,2]).
```

specifies that the lookup order for `[[]]` is a) the constructor name (`./2`), b) the first argument (not pre-indexed), and c) the second argument (pre-indexed). The resulting trie is in Fig. 2. In the figure, each node number represents a position in the trie. Singly circled nodes are temporary nodes, doubly circled nodes are final nodes. Final nodes encode terms. The initial node (```) is unique for each `rtype`. Labels between nodes indicate the lookup input. They can be constructor names (e.g., `./2`), `nv` terms (e.g., `b(y)`), or other pre-indexed `lst` (e.g., `#2` for `[], or `#5(z)` for `[c(z)]`). The arguments are placeholders for the non-indexed information. That is, a term `[a(g),c(h)]` would be encoded as `#9(g,h)`.

Trie indexing also supports anchoring on non-root nodes. Consider this declaration:

```prolog
:- rtype lst ---> []; [nv|lst]:.:index([2,0,1]).
```

Figure 3 shows the resulting trie (which has been separated into different subtrees for the sake of clarity). The lookup now starts from the second argument, then...
the constructor name, and finally the first argument. The main difference w.r.t. the previous indexing method is that the beginning node is another pre-indexed term. This may lead to more optimal memory layouts and need fewer lookup operations. Note that constructor names in the edges from initial nodes need to be prefixed with the name of the \( rtype \). This is necessary to avoid ambiguities, since the initial node is no longer unique.

**Garbage Collection and Indexing Methods.** Indexing structures require special treatment for garbage collection. In principle, it would not be necessary to keep in a trie nodes for terms that are no longer reachable (e.g., from the heap, WAM registers, or dynamic predicates), except for caching to speed-up node creation. Node removal may make use of lookup order. That is, if a key at a temporary level \( n \) corresponds to an atom that is no longer reachable, then all nodes above \( n \) can be safely discarded.

Anchoring on non-root nodes allows the simulation of interesting memory layouts. For example, a simple way to encode objects in Prolog is by introducing a

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6 Automatic garbage collection of indexing structures is not supported in the current implementation.
new object operation that creates new fresh atoms, and storing object attributes with a dynamic `objattr(ObjId, AttrName, AttrValue)` predicate. Anchoring on `ObjId` allows fast deletion (at the implementation level) of all attributes of a specific object when it becomes unreachable.

4 Applications and Experimental Evaluation

To show the feasibility of the approach, we have implemented the pre-indexing transformations as source-to-source transformations within the Ciao system. This is done within a Ciao package which defines the syntax and processes the `rtype` declarations as well as the marking of pre-indexing points.

As examples, we show algorithmically efficient implementations of the Lempel-Ziv-Welch (LZW) lossless data compression algorithm and the Floyd-Warshall algorithm for finding the shortest paths in a weighted graph, as well as some considerations regarding supporting module system implementation. In the following code, `forall/2` is defined as \( + (\text{Cond}, + \text{Goal}) \).

4.1 Lempel-Ziv-Welch compression

Lempel-Ziv-Welch (LZW) [16] is a lossless data compression algorithm. It encodes an input string by building an indexed dictionary \( D \) of words and writing a list of dictionary indices, as follows:

1- \( D := \{ w \mid w \text{ has length 1} \} \) (all strings of length one).
2- Remove from input the longest prefix that matches some word \( W \) in \( D \), and emit its dictionary index.
3- Read new character \( C \), \( D := D \cup \text{concat}(W, C) \), go to step 2; otherwise, stop.

A simple Prolog implementation is shown in Fig. 4 and Fig. 5. Our implementation uses a dynamic predicate `dict/2` to store words and corresponding numeric indices (for output). Step 1 is implemented in the `build_dict/1` predicate. Steps 2 and 3 are implemented in the `compress/3` predicate. For encoding words we use lists. We are only interested in adding new characters and word matching. For that, list construction and unification are good enough. We keep words in reverse order so that appending a character is done in constant time. For constant-time matching, we use an `rtype` for pre-indexing lists. The implementation is straightforward. Note that we add a character to a word in \( WC \)

---

7 We use updates in the dynamic program database as an instrumental example for showing the benefits of preindexing from an operational point of view. It is well known that this style of programming is often not desirable. The illustrated benefits of preindexing can be easily translated to more declarative styles (like declaring and composing effects in the type system) or more elaborate evaluation strategies (such as tabling, that uses memoization techniques).
Fig. 4. LZW Compression: Main code.

<table>
<thead>
<tr>
<th>data size</th>
<th>indexing (time)</th>
</tr>
</thead>
<tbody>
<tr>
<td></td>
<td>original</td>
</tr>
<tr>
<td>data1</td>
<td>1326</td>
</tr>
<tr>
<td>data2</td>
<td>83101</td>
</tr>
<tr>
<td>data3</td>
<td>149117</td>
</tr>
</tbody>
</table>

Table 1. Performance of LZW compression (in seconds) by indexing method.

= #lst([C\^W]) (Line 8). The annotation (whose syntax is implemented as a user-defined Prolog operator) is used by the compiler to generate the pre-indexed version of term construction. In this case, it indicates that words are pre-indexed using the \lst\ rtype and that \W is already pre-indexed (indicated by the escape ^ prefix). Thus we can effectively obtain optimal algorithmic complexity.

Performance evaluation. We have encoded three files of different format and size (two HTML files and a Ciao bytecode object) and measured the performance of alternative indexing and pre-indexing options. The experimental results for the algorithm implementation are shown in Table 1.\(^8\) The columns under indexing show the execution time in seconds for different indexing methods: none indicates that no indexing is used (except for the default first argument, first

\(^8\) Despite the simplicity of the implementation, we obtain compression rates similar to gzip.
% Mapping between words and dictionary index
:- data dict/2.
% NOTE: #lst can be changed or removed, ^ escapes cast
% Anchors to 2nd arg in constructor
:- rtype lst ---> []; [int|lst]:::index([2,0,1]).

build_dict(Size) :- % Initial dictionary
    assertz(dictsize(Size)),
    Size1 is Size - 1,
    forall(between(0, Size1, I), % Single code entry for I
        assertz(dict(#lst([I]), I))).
insert(W) :- % Add W to the dictionary
    retract(dictsize(Size)), Size1 is Size + 1, assertz(dictsize(Size1)),
    assertz(dict(W, Size)).

Fig. 5. LZW Compression: Auxiliary code and rtype definition for words.

level indexing); clause performs multi-level indexing on dict/2; term uses pre-indexed terms.

Clearly, disabling indexing performs badly as the number of entries in the dictionary grows, since it requires one linear (w.r.t. the dictionary size) lookup operation for each input code. Clause indexing reduces lookup complexity and shows a much improved performance. Still, the cost has a linear factor w.r.t. the word size. Term pre-indexing is the faster implementation, since the linear factor has disappeared (each word is uniquely represented by a trie node).

4.2 Floyd-Warshall

The Floyd-Warshall algorithm computes the shortest paths problem in a weighted graph in $O(n^3)$ time, where $n$ is the number of vertices. Let $G = (V, E)$ be a weighted directed graph, $V = v_1, \ldots, v_n$ the set of vertices, $E \subseteq V^2$, and $w_{i,j}$ the weight associated to edge $(v_i, v_j)$ (where $w_{i,j} = \infty$ if $(v_i, v_j) \notin E$ and $w_{i,i} = 0$). The algorithm is based on incrementally updating an estimate on the shortest path between each pair of vertices until the result is optimal. Figure 6 shows a simple Prolog implementation. The code uses a dynamic predicate dist/3 to store the computed minimal distance between each pair of vertices. For each vertex $k$, the distance between each $(i, j)$ is updated with the minimum distance calculated so far.

Performance evaluation. The performance of our Floyd-Warshall implementation for different sizes of graphs is shown in Fig. 7. We consider three indexing methods for the dist/3 predicate: def uses the default first order argument
floyd_warshall :-
% Initialize distance between all vertices to infinity
forall((vertex(I), vertex(J)), assertz(dist(I,J,1000000))),
% Set the distance from V to V to 0
forall(vertex(V), set_dist(V,V,0)),
forall(weight(U,V,W), set_dist(U,V,W)),
forall((vertex(K), vertex(I), vertex(J)),
  (dist(I,K,D1),
   dist(K,J,D2),
   D12 is D1 + D2,
   mindist(I,J,D12))).

mindist(I,J,D) :- dist(I,J,OldD), ( D < OldD -> set_dist(I,J,D) ; true ).

set_dist(U,V,W) :- retract(dist(U,V,_)), assertz(dist(U,V,W)).

Fig. 6. Floyd-Warshall Code

indexing, \( t12 \) computes the vertex pair key using two-level indices, \( p12 \) uses a packed integer representation (obtaining a single integer representation for the pair of vertices, which is used as key), and \( p12a \) combines \( p12 \) with a specialized array to store the \( \text{dist/3} \) clauses.

The execution times are consistent with the expected algorithmic complexity, except for \( \text{def} \). The linear relative factor with the rest of methods indicates that the complexity without proper indexing is \( O(n^4) \). On the other hand, the plots also show that specialized computation of keys and data storage (\( p12 \) and \( p12a \)) outperforms more generic encoding solutions (\( t12 \)).

4.3 Module System Implementations

Module systems add the notion of modules (as separate namespaces) to predicates or terms, together with visibility and encapsulation rules. This adds a significantly complex layer on top of the program database (whether implemented in C or in Prolog meta-logic as hidden tables, as in Ciao [4]). Nevertheless, almost no changes are required in the underlying emulator machinery or program semantics. Modular terms and goals can be perfectly represented as \( \text{M:T} \) terms and a program transformation can systematically introduce \( \text{M} \) from the context. However, this would include a noticeable overhead. To solve this issue, Ciao reserves special atom names for module-qualified terms (currently, only predicates).

We can see this optimization as a particular case of pre-indexing, where the last step in module resolution (which maps to the internal representation) is a pre-indexing cast for an \( \text{mpred rtype} \):

\[
: - \text{rtype mpred ---> nv:nv ::: index([1,0,2]).}
\]
Fig. 7. Execution time for Floyd-Warshall

For example, given a module $M = \text{lists}$ and goal $G = \text{append}(X,Y,Z)$, the pre-indexed term $MG = \#\text{mpred}(M;G)$ can be represented as $'\text{lists:append}'(X,Y,Z)$,\(^9\) where the first functor encodes both the module and the predicate name. To enable meta-programming, when $MG$ is provided, both $M$ and $G$ can be recovered.

Internally, another rewrite step replaces predicate symbols by actual pointers in the bytecode, which removes yet another indirection step. This indicates that it would be simple to reuse pre-indexing machinery for module system implementations, e.g., to enhance modules with hierarchies or provide better tools for meta-programming. In principle, pre-indexing would bring the advantages of efficient low-level code with the flexibility of Prolog-level meta representation of modules. Moreover, anchoring on $M$ mimicks a memory layout where predicate tables are stored as key-value tables inside module data structures.

5 Related Work

There has been much previous work on improving indexing for Prolog and logic programming. Certain applications involving large data sets need any- and multi-argument indexing. In [7] an alternative to static generation of multi-argument indexing is presented. The approach presented uses dynamic schemes for demand-driven indexing of Prolog clauses. In [14] a new extension to Prolog indexing is proposed. User-defined indexing allows the programmer to index both instantiated and constrained variables. It is used for range queries and spatial queries, and allows orders of magnitude speedups on non-trivial datasets.

Also related is ground-hashing for tabling, studied in [17]. This technique avoids storing the same ground term more than once in the table area, based on

\(^9\) Note that the identifier does not need any symbolic description in practice.
computation of hash codes. The approach proposed adds an extra cell to every compound term to memoize the hash code and avoid the extra linear time factor.

Our work relates indexing techniques (which deal with fast lookup of terms in collections) with term representation and encoding (which clearly benefits from specialization). Both problems are related with optimal data structure implementation. Prolog code is very often used for prototyping and then translated to (low-level) imperative languages (such as C or C++) if scalability problems arise. This is however a symptom that the emulator and runtime are using suboptimal data structures which add unnecessary complexity factors. Many specialized data structures exist in the literature, with no clear winner in all cases. If they can be directly implemented in Prolog, they are often less efficient than their low-level counterparts (e.g., due to data immutability). Without proper abstraction they obscure the program to the point where a low-level implementation may not be more complex. On the other hand, adding them to the underlying Prolog machines is not trivial. Even supporting more than one term representation may have prohibitive costs (e.g., efficient implementations require a low number of tags, small code that fits in the instruction cache, etc.). Our work aims at reusing the indexing machinery when possible and specializing indexing for particular programs.

The need for the right indexing data structures to get optimal complexity is also discussed in [11] in the context of CHR. In [9], an improved term encoding for indexed ground terms that avoids the costs of additional hash-tables is presented. This has similar results to anchoring in pre-indexing. Reusing the indexing machinery is also studied in [8], which shows term flattening and specialization transformations.

6 Conclusions and Future Work

Traditionally, Prolog systems index terms during clause selection (in the best case, reducing a linear search to constant time). Despite that, index lookup is proportional to the size of the term. In this paper we have proposed a mixed approach where indexing is precomputed during term creation. To do that, we define a notion of instantiation types and annotated constructors that specify the indexing mode. The advantage of this approach is that lookups become sub-linear. We have shown experimentally that this approach improves clause indexing and that it has other applications, for example for module system implementation.

These results suggest that it may be interesting to explore lower-level indexing primitives beyond clause indexing. This work is also connected with structure sharing. In general, pre-indexing annotations allow the optimization of simple Prolog programs with scalability problems due to data representation.

As future work, there are some open lines. First, we plan to polish the current implementation, which is mostly based on program rewriting and lacks garbage
collection of indexing tables. We expect major performance gains by optimizing some operations at the WAM or C level. Second, we want to extend our repertoire of indexing methods and supporting data structures. Finally, rtype declarations and annotations could be discovered and introduced automatically via program analysis or profiling (with heuristics based on cost models).

References