Abstract Program Slicing: An Abstract Interpretation-Based Approach to Program Slicing

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In the present article, we formally define the notion of abstract program slicing, a general form of program slicing where properties of data are considered instead of their exact value. This approach is applied to a language with numeric and reference values and relies on the notion of abstract dependencies between program statements.

The different forms of (backward) abstract slicing are added to an existing formal framework where traditional, nonabstract forms of slicing could be compared. The extended framework allows us to appreciate that abstract slicing is a generalization of traditional slicing, since each form of traditional slicing (dealing with syntactic dependencies) is generalized by a semantic (nonabstract) form of slicing, which is actually equivalent to an abstract form where the identity abstraction is performed on data.

Sound algorithms for computing abstract dependencies and a systematic characterization of program slices are provided, which rely on the notion of agreement between program states.

CCS Concepts:
- Security and privacy → Logic and verification;
- Theory of computation → Logic and verification; Program verification; Program analysis; Automated reasoning; Programming logic; Hoare logic; Pre- and post-conditions;
- Software and its engineering → Formal software verification;

Additional Key Words and Phrases: Program slicing, semantics, static analysis, abstract interpretation

ACM Reference Format:
Isabella Mastroeni and Damiano Zanardini. 2017. Abstract program slicing: An abstract interpretation-based approach to program slicing. ACM Trans. Comput. Logic 18, 1, Article 7 (February 2017), 58 pages.
DOI: http://dx.doi.org/10.1145/3029052

1. INTRODUCTION

It is well known that, as the size of programs increases, it becomes impractical to maintain them as monolithic structures. Indeed, splitting programs into smaller pieces allows one to construct, understand, and maintain large programs much more easily. Program slicing [Weiser 1984; Tip 1995; Binkley and Gallagher 1996; De Lucia 2001] is a program manipulation technique that extracts, from programs, those statements that are relevant to a particular computation. In the most traditional definition, a program slice is an executable program whose behavior must be identical to a specific subset of the original program’s behavior. The specification of this subset is called the slicing criterion and can be expressed as the value of some set of variables at some...
Fig. 1. $Q$, $R$, and $S$ are, respectively, a slice, a semantic slice, and an abstract slice of $P$.

set of statements and/or program points [Weiser 1984]. Slicing$^1$ can be and is used in several areas like debugging [Weiser 1984], software maintenance [Gallagher and Lyle 1991], comprehension [Canfora et al. 1998; Field et al. 1995], and re-engineering [Cimitile et al. 1996].

Since the seminal paper introducing slicing [Weiser 1984], there have been many works proposing several notions of slicing and different algorithms to compute slices (see Tip [1995] and De Lucia [2001] for good surveys about existing slicing techniques). Program slicing is a transformation technique that reduces the size of programs to analyze. Nevertheless, the reduction obtained by means of standard slicing techniques may be not sufficient for simplifying program analyses since it may keep more statements than those strictly necessary for the desired analysis. Suppose we are analyzing a program, and suppose we want a variable $x$ to have a particular property $\rho$ at a given program point $n$. If we realize that $x$ does not have that property at $n$, then we may want to understand which statements affect that property of $x$ in order to find out more easily where the computation went wrong. In this case, we are not interested in the exact value of $x$, so that we may not need all the statements that a standard slicing algorithm would return. Instead, we would need a technique that returns the minimal amount of statements that actually affect that specific property of $x$.

1.1. Abstract Program Slicing

This article introduces and discusses a “semantic” general notion of slicing, called abstract program slicing, looking for those statements affecting a property (modeled in the context of abstract interpretation [Cousot and Cousot 1977]) of a set of variables of interest, the so-called abstract criterion. The idea behind this new notion of slicing is investigating more semantically precise notions of dependency between variables. In other words, when a syntactic dependency is detected, such as the dependency, in an assignment, of defined variables from used variables, we look further for semantic dependencies, that is, dependencies between values of variables.

Consider the program $P$ in Figure 1, and suppose that we are interested in the variable $d$ at the end of the execution. Standard slicing algorithms extract slices by computing syntactic dependencies; in this sense, $d$ depends on all $c$, $b$, and $a$, so that a sound slice would have to take all the statements involving all these variables. In the figure, $Q$ is a slice of $P$ with respect to that criterion. However, if we are interested in a more precise, semantic notion of slicing, then we could observe that the value of $d$ only depends on the values of variables $c$ and $b$, so that a more precise slice would be represented by $R$. Finally, if we are interested in the parity of $d$ at that point, then we observe that parity of $d$ does not depend on the value of $c$, and $S$ is an abstract slice of $P$ with respect to the specified criterion. Even in this simple case, the abstract slice gives more precise information about the statements affecting the property of interest.

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1We use slicing (slice) and program slicing (program slice) as interchangeable terms.
1.2. Outline and Contributions

In this article, we aim at introducing a generalized notion of slicing, allowing us to weaken the notion of “dependency” (from syntax, to semantics, to abstract semantics) with respect to what is considered relevant for computing the slice. Since our generalization is a semantic one, we start from the unifying framework by Binkley et al. [2006a, 2006b], recalled and discussed in Section 3 and the appendix, where different forms of slicing are defined and compared w.r.t. their characteristics (static/dynamic, iteration-count/non-iteration-count, etc.) into a comprehensive formal framework. The structure of this framework is based on the formal definition of the criterion, inducing a semantic equivalence relation $E$ that uniquely characterizes the set of possible slices of a program $P$ as the set of all the subprograms equivalent to $P$ with respect to $E$.

This structure makes the framework suitable for the introduction and the formal definition of an abstract form of slicing in Section 4, since abstraction corresponds simply to consider a weaker criterion, which implies weakening the equivalence relation $E$ defining slicing.

Once we have the equivalence relation defining a desired notion of slicing w.r.t. a given criterion, we show in Section 5 how this corresponds to fixing the notion of dependency we are interested in (namely, the notion of dependency determining what has to be considered relevant in the construction of slicing), and we show how the extension to semantic dependencies may be used to extend the program-dependency-graph-based approach to computing slices [Horwitz et al. 1989]. Finally, we define in Section 6 a notion of abstract dependencies implying abstract criteria. In Section 6.1, we show that this new notion of dependency is not suitable for computing slices by using program dependency graphs, and in Section 6.2 propose the $\text{findNdeps}$ algorithm for computing (abstract) dependencies and the $E_{\text{dep}}$ algorithm for finding an abstract domain for which some specific abstract dependency of an expression on some variables does not hold.

Next, Section 7 gives a systematic approach to compute backward slices. Such an approach relies on two systems of logical rules in order to prove (1) Hoare-style tuples capturing the effect of executing a statement $s$ on a pair of states for which some similarity (agreement) is required by the slicing criterion (indeed, this similarity corresponds to the semantic equivalent relation) and (2) when some properties of the state do not change (are preserved) after $s$ is executed. The combination of the results provided by these rule systems allows one to decide whether it is safe to remove a statement from a program without changing the observation corresponding to the criterion. Importantly, the rule systems and algorithms provided in Section 7 rely on the knowledge and manipulation of a “library” of abstract properties. For example, in order to infer that $2 \times x$ is always even, the abstract domain representing the parity of number must be known and available as a “component” of the logical systems. If no abstract property is known except the identity (which is the most precise property and is not really abstract), then the approach boils down to standard slicing. Importantly, it becomes clear in this case that slices on the same variables (properties of them in the abstract case; exact values in the concrete case) are generally bigger in the concrete setting (when identity is the only available property) than in the corresponding abstract slicing. Needless to say, this does not mean that every algorithm for abstract slicing will perform better than any algorithm for nonabstract slicing; rather, it provides a practical insight of how optimal (purely semantic-based) abstract slices may not include statements that are included in concrete slices.

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2The framework proposed in Binkley et al. [2006a, 2006b] is parametric on the syntactic relation, but here we only consider the relation of being a subprogram.
Part of this work has been previously published in conference proceedings [Mastroeni and Zanardini 2008; Zanardini 2008; Mastroeni and Nikolić 2010]. The present article joins these works into a coherent framework and contains a number of novel contributions:

- In Section 4, we formally prove that abstract slicing can be put in the formal framework of Binkley et al. [2006a] and generalizes concrete forms of slicing.
- In the same section, we formally define the notion of dependency induced by a particular criterion, that is, by the equivalence relation among programs corresponding to the chosen criterion.
- In Section 5, we define and prove how we can approximate this (concrete semantic) dependency in order to use it for pruning PDGs and computing slicing with the well-known PDG-based algorithm for slicing [Reps and Yang 1989].
- In Section 6.1, we discuss why the idea of pruning PDGs is not applicable to the abstract notion of dependency, thus showing the need of providing different approaches for computing abstract slices.
- The discussion on the findNdeps and Edep algorithm in Sections 6.2.1 and 6.2.2 has been significantly improved, and several examples have been provided.
- The treatment of nonnumerical values when computing slices was already considered in Zanardini [2008]. However, the language under study in the present article is different in that it is closer to standard object-oriented languages. More concretely, that work used complex identifiers \(x.f.g\) as if they were normal variables, thus obtaining that sharing between variables was easier to deal with. However, this came at the cost of increasing the number of “variables” to be tracked by the analysis. Section 7 discusses the logical machinery to compute backward slices in this context. Moreover, examples have been provided to illustrate how properties of the heap can be taken into account.
- The \(G\)-system introduced in Section 7.1 is a quite refined version of the \(A\)-system [Zanardini 2008]; rules for variable assignment and field update have been changed according to the new language (which implies a number of technical issues); there is a new rule \(G\)–ID; the overall discussion has been improved.
- The rule system for proving the preservation of properties (the \(PP\)-system) is explicitly introduced in Section 7.1.7.
- The description of how statements can be erased has been greatly improved; an algorithm has been explicitly introduced in Section 7.3, which labels each program point with agreements according to the \(G\)-system. A thorough discussion and proofs are provided, so that it is guaranteed that the conditions for erasing a statement (relying on the \(G\)-system, the \(PP\)-system, and the labelSequence procedure for labeling program points with agreements) are sound.
- In Section 7.4.2, recent work on field-sensitive sharing analysis [Zanardini 2015] is included in the computation of abstract slices, which results in improving the precision when data structures in the heap overlap.

2. PRELIMINARIES

2.1. The Programming Language

The language is a simple imperative language with basic object-oriented features, whose syntax will be easy to understand for anyone who is familiar with imperative programming and object orientation. The language syntax includes the usual arithmetic expressions \(\text{Exp}\) and access to object fields via “dot” selectors. A statement can be \textbf{skip}, a variable assignment \(x:=e\), a field update \(x.f:=e\), a conditional, or a \textbf{while} loop. In addition, there exist special statements: (1) \textbf{read}, which reads the value of some variable from the input, simulating the use of parameters—this kind of statement can
only appear at the beginning of the program, and (2) write, which can only appear at
the end of the program and outputs the current value of some variables.\textsuperscript{3} For simplicity,
guards in conditionals and loops are supposed not to have \textit{side effects}. We denote by $\wp$ the
set of all programs.

$X$ is the set of program variables and $V$ denotes the set of values, which can be either
integer or reference values, or the \texttt{null} constant ($V = \mathbb{Z} \cup \mathbb{R} \cup \{\texttt{null}\}$); every variable is
supposed to be well typed (as integer or reference) at every program point. $L$ denotes
the set of \textit{line numbers} (program points). Let $l \in L$, and $\text{Sim}(l)$ be the statement at
program line $l$. For a given program $P$, we denote by $L_P \subseteq L$ the set of all and only
the line numbers corresponding to statements of the program $P$, that is, $L_P = \{ l \in L \mid \text{Sim}(l) \in P \}$. This definition is necessary since when we look for slicing, we erase
statements without changing the numeration of line numbers; for instance, in Figure 1,
we have that $\text{Sim}(4) \notin Q$, so that $L_Q = \{1, 2, 3, 5\}$.

A \textit{program state} $\sigma \in \Sigma$ is a pair $\langle n^k, \mu \rangle$, where $n$ is the executed program point, $k$ is
the number of times the statement at $n$ has been reached so far, and $\mu$ is the memory.
A \textit{memory} is a pair $\langle \varepsilon, h \rangle$, where the \textit{store} $\varepsilon : X \rightarrow V$ maps variables to values, and
the \textit{heap} $h$ is a sequence of locations where objects can be stored; a reference value
corresponds to one of such locations. An \textit{object} $o$ maps field identifiers to values, in
the usual way; $o.f$ is the value corresponding to the field $f$ of the object $o$ and can be
either a number, the location in which another object is stored, or \texttt{null}. For the sake
of simplicity, \textit{classes} are supposed to be declared somewhere, and field accesses are
supposed to be consistent with class declarations.

Unless ambiguity may arise, a memory (or even an entire program state) can be
considered directly as a store, so that $\mu(x) (\sigma(x), \text{respectively})$ will be the value of $x$
in the store contained in $\mu$ (in $\sigma$, respectively). Moreover, a store $\varepsilon$ can be represented as $\{x_1 \leftarrow v_1, \ldots, x_m \leftarrow v_m\}$, meaning that $\varepsilon(x_i) = v_i$ for every $i$, and again, $\mu = \{x_1 \leftarrow v_1, \ldots, x_m \leftarrow v_m\}$ (respectively) can be used instead of $\varepsilon = \{x_1 \leftarrow v_1, \ldots, x_m \leftarrow v_m\}$ whenever the store is the only relevant part of the memory
(the state, respectively).

A \textit{state trajectory} $\tau \in \mathbb{T} = \Sigma^*$ is a sequence of program states through which a
program goes during the execution. State trajectories are actually traces equipped with
the $\bar{k}$ component. The state trajectory obtained by executing program $P$ from
the input memory $\mu$ is denoted $\tau^P_\mu$. Moreover, $\tau[n]$ will be the set of states in $\tau$ where
the program point is $n$. Any initial state has $n = 1$; that is, the set of initial states is
$\Sigma_i = \{ \langle 1^1, \mu \rangle \mid \mu \in M \}$.

In the following, $[\cdot] : P \times \wp(\Sigma_i) \rightarrow \wp(\mathbb{T})$ denotes the program semantics where $[P](S)$
returns the set of state trajectories obtained by executing the program $P$ starting from
any initial state in $S \subseteq \Sigma_i$, that is, $[P](S) = \{ \tau^P_\mu \mid \langle 1^1, \mu \rangle \in S \}$. We abuse notation by
denoting in the same way also the semantics of expressions, namely, $[\cdot] : \wp \times \Sigma \rightarrow V$,
which is such that $[e](\sigma) (e \in \wp)$ returns the evaluation of $e$ in $\sigma$. Finally, if $S \subseteq \Sigma$
for the sake of simplicity, we still abuse notation by denoting in the same way also the
additive lift of semantics, that is, $[e](S) = \{ [e](\sigma) \mid \sigma \in S \}$.

### 2.2. Basic Abstract Interpretation

This section introduces the lattice of \textit{abstract interpretations} [Cousot and Cousot 1977].
Let $\langle C, \leq, \vee, \wedge, \top, \bot \rangle$ denote a complete lattice $C$, with ordering $\leq$, lub $\vee$, glb $\wedge$, and
top and bottom element $\top$ and $\bot$, respectively. A \textit{Galois connection} (G.c.) is a
pair of monotone functions $\alpha : C \rightarrow A$ and $\gamma : A \rightarrow C$ such that $\alpha(x) \leq_A y \iff x \leq_C \gamma(y)$. In standard terminology, $C$ and $A$ are, respectively, the concrete and the

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\textsuperscript{3}As a matter of fact, this kind of statement is only included for back-compatibility and readability.
abstract domain. Abstract domains can be formulated as upper closure operators ($\rho$) [Cousot and Cousot 1977]. Given an ordered set $C$ with ordering $\leq_C$, a uco on $C$, $\rho : C \rightarrow C$, is a monotone, idempotent ($\rho(\rho(x)) = \rho(x)$), and extensive ($\forall x \in C. \ x \leq_C \rho(x)$) map. Each uco $\rho$ is uniquely determined by the set of its fix-points, which is its image; that is, $\rho(C) = \{ x \in C \mid \rho(x) = x \}$. When $C = \wp(D)$ for some set $D$, and $v \in D$, then we usually write $\rho(v)$ instead of $\rho(\{v\})$ (and in general for any function, $f(v)$ instead of $f(\{v\})$). If $C$ is a complete lattice, then $\langle uco(C), \subseteq, \cup, \cap, \land, x. T, \land, x.x \rangle$ is a complete lattice, where $uco(C)$ is the domain of all the upper closure operators on the lattice $C$; for every two ucoss $\rho_1, \rho_2 \in uco(C)$, $\rho_1 \subseteq \rho_2$ if and only if $\forall y \in C. \ \rho_1(y) \leq \rho_2(y)$ if and only if $\rho_2(C) \subseteq \rho_1(C)$; and, for every $\{\rho_i\}_{i \in I} \subseteq uco(C)$, $(\cap_{i \in I} \rho_i)(x) = \land_{i \in I} \rho_i(x)$ and $(\cup_{i \in I} \rho_i)(x) = x \Leftrightarrow \forall i \in I. \ \rho_i(x) = x$. In the following, we will denote by $\rho_{\text{AB}}$ the most concrete uco on a domain, that is, $\lambda x.x$, and by $\rho_{\text{T}}$ the most abstract one $\land x.x$. $A_1$ is more precise than $A_2$ (i.e., $A_2$ is an abstraction of $A_1$) if and only if $A_1 \subseteq A_2$ in $uco(C)$. The reduced product of a family $\{\rho_i\}_{i \in I}$ is $\cap_{i \in I} \rho_i$ and is one of the best-known operations for composing domains.

Example 2.1 (Numerical abstract domains). Let the concrete domain $C$ be $\wp(\mathbb{Z})$; the parity abstract domain $\rho_{\text{Par}}$ in Figure 2 (on the left) represents the parity of numbers and is determined by fix-points $\{\text{bot}, \text{even}, \text{odd}, \text{top}\}$, where $\text{even}$ and $\text{odd}$ denote even and odd numbers, respectively; $\text{bot}$ is the empty set; and $\text{top} = \mathbb{Z}$. For example, $\rho_{\text{Par}}(\{2, 4, 10\}) = \text{even}$ (all numbers are even), $\rho_{\text{Par}}(\{3, 7\}) = \text{odd}$ (both numbers are odd), and $\rho_{\text{Par}}(\{4, 5\}) = \text{top}$ (there are both even and odd numbers).

The sign abstract domain $\rho_{\text{Sign}}$ (in the center of Figure 2) is characterized by fix-points $\{\text{bot}, \text{zero}, \text{pos}, \text{neg}, \text{top}\}$ and tracks the sign of integers (zero, positive, negative, etc.). For example, $\rho_{\text{Sign}}(\{0\}) = \text{zero}$, $\rho_{\text{Sign}}(\{-3, -4, -5\}) = \text{neg}$, $\rho_{\text{Sign}}(\{1, 2, 4\}) = \text{pos}$, and $\rho_{\text{Sign}}(\{1, -1\}) = \text{top}$. Finally, the parity-sign domain $\rho_{\text{ParSign}}$ (on the right), which is the reduced product $\cap$ of $\rho_{\text{Par}}$ and $\rho_{\text{Sign}}$, captures both properties (the parity and the sign) and has fix-points $\{\text{bot}, \text{zero}, \text{pos}, \text{neg}, \text{top}\}$. Formally speaking, the value of a reference variable is either a location $\ell$ or $\text{null}$. However, the domains introduced in the next example classify variables not only with respect to $\ell$ itself but also on the data structure in the heap, which is reachable from $\ell$. This point of view is similar to previous work on static analysis of properties of the heap like sharing [Secchi and Spoto 2005] or cyclicity [Rossignoli and Spoto 2006; Genaim and Zanardini 2013].

Example 2.2 (Reference abstract domains). Let $C$ be $\wp(\mathbb{R} \cup \{\text{null}\})$, that is, the possible values of reference variables. The nullity domain $\rho_{\text{null}}$ classifies values on
nullity and has fix-points \([\text{bot}], [\text{null}], [\text{non–null}], [\text{top}]\), where the concretizations of \([\text{null}]\) and \([\text{non–null}]\) are, respectively, \([\text{null}]\) and \([\mathbb{R}]\).

On the other hand, it is possible to define a cyclicity domain \(\rho_{\text{CYC}}\) that classifies variables on whether they point to cyclic or acyclic data structures [Genaim and Zanardini 2013]. A cycle in the heap is a path in which the same location is reached more than once; a double-linked list (one that can be traversed in both directions) is a good example of a cyclic data structure. The fix-points of this domain are \([\text{bot}], [\text{cyc}], [\text{acyc}], [\text{top}]\), where all acyclic values (including \([\text{null}]\)) are abstracted to \([\text{acyc}]\), and all cyclic values (i.e., locations from which a cycle is reachable) are abstracted to \([\text{cyc}]\). Both domains and their reduced product are depicted in Figure 3; note that there are no values that are both null and cyclic, so that their intersection collapses to \([\text{bot}]\).

Finally, the identity domain \(\rho_{\text{id}}\) abstracts two concrete values to the same abstract value only if they are equal. Two references are equal if (1) they are both null or (2) they are both nonnull and the objects stored in the corresponding locations are equal, where equality on objects means that all their numeric fields must be the same number and all reference fields must be equal (w.r.t. this same notion of equality on references).

Let us consider now \(D = C^n\) (\(C\) lattice and \(n \in \mathbb{N}\)), namely, \(x \in D\) is an \(n\)-tuple of elements of \(C\), and consider \(\rho \in uco(D)\). In this case, we can distinguish between two kinds of abstractions: nonrelational and relational abstractions [Cousot 2001; Cousot and Cousot 1979]. The nonrelational or attribute-independent one [Cousot and Cousot 1979, Example 6.2.0.2] consists of ignoring the possible relationships between the values of the abstracted inputs. For instance, if \(\rho\) is applied to the values of variables \(x\) and \(y\), then \(\rho\) can be approximated through projection by a pair of abstractions on the single variables, analyzing the single variables in isolation. For the sake of simplicity, without losing generality, consider \(n = 2\), that is, \(D = C^2 = C \times C\). Formally, \(\rho \in uco(C \times C)\) is nonrelational if there exist \(\delta_1, \delta_2 \in uco(C)\) such that \(\rho(x, y) = (\delta_1(x), \delta_2(y))\), that is, \(\rho \in uco(C) \times uco(C) \subseteq uco(C \times C)\). For instance, let \(\rho_{\text{PAR}}\) be the abstract domain depicted in Figure 2 expressing the parity of integer values; the \(\rho_{\text{PAR}}\) nonrelational property of \((x, y)\) provides the parity of \(x\) and \(y\) independently one from each other, meaning that all the possible combinations of parity of \(x\) and \(y\) are possible as results \(\{(\text{even},\text{even}),(\text{even},\text{odd}),(\text{odd},\text{even}),(\text{odd},\text{odd})\}\) and all combinations where at least one variable is [top] or [bot]). Relational abstractions may preserve some of the relationship between the analyzed values [Cousot 2001]. For instance, we could define an abstraction preserving the fact that \(x\) is even ([even]) if and only if \(y\) is odd ([odd]). It is clear that, in this case, we are more precise since the only possible analysis results are \(\{(\text{even},\text{odd}),(\text{odd},\text{even}),(\text{top},\text{top}),(\text{bot},\text{bot})\}\).

If \(\rho \in uco(C)\), \(f \in C \longrightarrow C\), and \(f^2 \in \rho(C) \longrightarrow \rho(C)\), then \(f^2\) is a sound approximation of \(f\) if \(\rho \circ f \subseteq f^2 \circ \rho\). \(f^\rho \triangleq \rho \circ f \circ \rho\) is known as the best correct approximation (BCA) of \(f\) in \(\rho\), which is always sound by construction. Soundness naturally implies fix-point
soundness, that is, $\rho(\text{lfp}_{\leq}^c f) \preceq \rho(\text{fixp})$. If $\rho \circ f = \rho \circ f \circ \rho$, then we say that $f^\rho$ is a complete approximation of $f$ [Cousot and Cousot 1979; Giacobazzi et al. 2000]. In this case, $\rho(\text{lfp}_{\leq}^c f) = \text{lfp}_{\leq}^c f^\rho$.

2.3. Equivalence Relations, Abstractions, and Partitions

Closure operators and equivalence relations are related concepts [Cousot and Cousot 1979]. Recently, this connection has been further studied in the field of abstract model checking and language-based security [Ranzato and Tapparo 2002; Hunt and Mastroeni 2005]. In particular, there exists an isomorphism between equivalence relations and a subclass of upper closure operators. Consider a set $S$: for each equivalence relation $\mathcal{R} \subseteq S \times S$, we can define an upper closure operator, $\text{Clo}^\mathcal{R} \in \text{uco}(\wp(S))$, such that $\forall x \in S \cdot \text{Clo}^\mathcal{R}(\{x\}) = [x]_\mathcal{R}$ and $\forall X \subseteq S \cdot \text{Clo}^\mathcal{R}(X) = \bigcup_{x \in X} [x]_\mathcal{R}$. Conversely, for each upper closure operator $\eta \in \text{uco}(\wp(S))$, we are able to define an equivalence relation $\text{Rel}^\eta \subseteq S \times S$ such that $\forall x, y \in S \cdot x \text{ Rel}^\eta y \iff \eta(x) = \eta(y)$. It is immediate to prove that $\text{Rel}^\eta$ is an equivalence relation, and this comes from $\eta$ being merely a function, not necessarily a closure operator. $\text{Clo}^\mathcal{R}$ is identified as the most concrete closure $\eta$ such that $\mathcal{R} = \text{Rel}^\eta$ [Hunt and Mastroeni 2005]. It is possible to associate with each upper closure operator the most concrete closure inducing the same partition on the concrete domain $S$:

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\Pi(\eta) \overset{\text{def}}{=} \text{Clo}^{\text{Rel}^\eta}.
$$

Note that, for all $\eta \in \text{uco}(\wp(S))$, $\Pi(\eta)$ is the (unique) most concrete closure that induces the same equivalence relation as $\eta$ ($\text{Rel}^\eta = \text{Rel}^{\Pi(\eta)}$). The fix-points of $\Pi$ are called the partitioning closures. With $\wp(S)$ being a complete Boolean lattice, an upper closure operator $\eta \in \text{uco}(\wp(S))$ is partitioning, that is, $\eta = \Pi(\eta)$, if and only if it is complemented, namely, if $\forall X \in \eta, \overline{X} \overset{\text{def}}{=} S \setminus X \in \eta$ [Hunt and Mastroeni 2005].

**Example 2.3.** Consider the set $S = \{1, 2, 3, 4\}$ and one of its possible partitions $\pi = \{\{1\}, \{2, 3\}, \{4\}\}$. The closure $\eta$ with fix-points $\emptyset, \{1\}, \{4\}, \{1, 2, 3\}, S$ induces exactly $\pi$ as a state partition, but the most concrete closure that induces $\pi$ is $\text{Clo}^\pi = \Pi(\eta) = \bigvee(\emptyset, \{1\}, \{2, 3\}, \{4\}, S)$, which is the closure on the right of Figure 4.

Given a partitioning upper closure operator $\rho$, an atom is an element $a$ of $\rho$ such that there does not exist another element $b$ with $\text{bot} \sqsubset b \sqsubset a$. For example, the atoms of $\rho_{\text{parSign}}$ are $\text{[poseven]}, \text{[posodd]}, \text{[zero]}, \text{[negeven]}, \text{[negodd]}$. In partitioning closures, atoms are all the possible abstractions of singletons: in fact, $\rho_{\text{parSign}}(\{n\})$ will never give $\text{[pos]}$ or $\text{[odd]}$ since there is always a more precise abstract value describing $n$. In the following, $\text{ATOM}_{\rho}(a)$ holds if and only if $a$ is an atom of $\rho$. 

ACM Transactions on Computational Logic, Vol. 18, No. 1, Article 7, Publication date: February 2017.
2.4. Abstract Semantics and Static Analysis

An abstract program semantics is the abstract counterpart of the concrete semantics w.r.t. an abstract program observation: it is meant to compute, for each program point, an abstract state that soundly represents invariant properties of variables at that point. In general, it is computed by an abstract interpreter [Cousot and Cousot 1979] collecting the set of all the possible values that each variable may have in each program point and abstracting this set in the chosen abstract domain.

Given a concrete program state $\sigma$ and an abstract domain $\rho \in \text{uco}(\varphi(\mathbb{V}))$, an abstract state $\sigma^\rho \in \Sigma^\rho$ is obtained by applying the abstraction $\rho$ to the values of variables stored in it. Namely, $\sigma^\rho = \langle \mathbb{V}^\rho, \mu^\rho \rangle$, where $\mu^\rho = \langle \varepsilon^\rho \rangle$ and $\varepsilon^\rho$ is such that $\varepsilon^\rho(x) = \rho(e(x))$. For simplicity, we can write $\sigma^\rho(x) = \rho(\sigma(x))$, treating the whole state as a store when applied to variables. In the case of a reference variable $x$, the abstraction $\sigma^\rho(x)$ gives information about the data structure pointed to by $x$ (e.g., if $\rho = \rho_{\text{CYC}}$, the cyclicity of the data structure can be represented). This explains why the heap is not represented explicitly in the abstract state: instead, relevant information about the heap is contained in the abstraction of variables (see the previous discussion before Example 2.2).

In the following, ordering $\leq$ on abstract states is variable-wise comparison between abstract values:

$$\sigma^\rho_1 \leq \sigma^\rho_2 \iff \forall x. \sigma^\rho_1(x) \subseteq \sigma^\rho_2(x).$$

The greater an abstract state is, the wider is the set of concrete states it represents; if $\sigma^\rho_1 \leq \sigma^\rho_2$, then $\sigma^\rho_1$ is called a refinement of $\sigma^\rho_2$. Moreover, a covering of $\sigma^\rho$ is a set of refinements $\{\sigma^\rho_1, \ldots, \sigma^\rho_n\}$ such that $\forall \sigma^\rho_i = \sigma^\rho$. The set of abstract state trajectories is $\mathcal{T}^\rho = \Sigma^{\rho^\ast}$; that is, an abstract trajectory is the computation of a program on the set of abstract states. The trace in $\mathcal{T}^\rho$ of a program $P$, starting from the abstract memory $\mu^\rho$, is denoted by $\tau^\rho_P$.

The abstract program semantics $\llbracket \cdot \rrbracket^\rho : \mathbb{P} \times \Sigma^\rho \rightarrow \mathcal{T}^\rho$ is such that $\llbracket P \rrbracket^\rho(S) = \{ \tau^\rho_P | (1^1, \mu^\rho) \in S \}$ is the set of the sequences of abstract states computed starting from the abstract initial states in $S \in \varphi(\Sigma^\rho)$. We also abuse notation by denoting $\llbracket \cdot \rrbracket^\rho$ also as the abstract evaluation of expressions. Namely, $\llbracket \cdot \rrbracket^\rho : \mathbb{E} \times \Sigma^\rho \rightarrow \rho(\mathbb{V})$ is such that $\forall x. \llbracket e \rrbracket^\rho(\sigma^\rho(x)) = \rho(\llbracket e \rrbracket(\sigma(x))) = \rho(e(\rho(\sigma(x))))$. This definition is correct, since by construction, we have that any abstract state $\sigma^\rho$ corresponds to a set of concrete states, that is, $\sigma^\rho = \{ \sigma \in \Sigma | \sigma = \sigma^\rho \} = \{ \sigma \in \Sigma | \forall x \in \mathbb{X}. \sigma^\rho(x) = \sigma(x) \}$; namely, it is the set of all the concrete states having as abstraction in $\rho$ precisely $\sigma^\rho$, and we abuse notation by denoting with $\llbracket \cdot \rrbracket^\rho$ also its additive lift. In other words, $\llbracket e \rrbracket^\rho$ is the best correct approximation of $[e]$ by means of an abstract value in $\rho$. In general, in order to compute the abstract semantics of a program on an abstract domain $\rho$, we have to equip the domain $\rho$ with the abstract versions of all the operators used for defining expressions. In our language, we should define, for example, the meaning of $+$, $-$, $\ast$, and $/$ on abstract values, that is, on sets of concrete values. This is standard in abstract interpretation, and these operations are defined for many known numerical abstract domains. For instance, the sound approximation of the sum operation on $\rho_{\text{PAR}}$ is the following:

$$\begin{align*}
\llbracket \text{even} \rrbracket + \llbracket \text{even} \rrbracket &= \llbracket \text{even} \rrbracket & \llbracket \text{top} \rrbracket + \bot &= \llbracket \text{top} \rrbracket \\
\llbracket \text{even} \rrbracket + \llbracket \text{odd} \rrbracket &= \llbracket \text{odd} \rrbracket & - + \llbracket \text{top} \rrbracket &= \llbracket \text{top} \rrbracket \\
\llbracket \text{odd} \rrbracket + \llbracket \text{even} \rrbracket &= \llbracket \text{odd} \rrbracket & \llbracket \text{bot} \rrbracket + \bot &= \llbracket \text{bot} \rrbracket \\
\llbracket \text{odd} \rrbracket + \llbracket \text{odd} \rrbracket &= \llbracket \text{even} \rrbracket & - + \llbracket \text{bot} \rrbracket &= \llbracket \text{bot} \rrbracket
\end{align*}$$

We can reason similarly for all the other operators.
In the following, the notation \([e]^{\rho}(S)\), where \(S\) is a set of concrete states, is used to denote the abstract evaluation of \(e\) in an abstract state, which is obtained by abstracting \(S\); as usual, \([e]^{\rho}(\sigma)\) stands for \([e]^{\rho}(\{\sigma\})\). The use of \([\cdot]^{\rho}\) in Section 6.2.1 and later in the article is twofold: (1) to infer invariant properties, as in Example 2.4, and (2) to evaluate expressions at the abstract level.

**Example 2.4.** Consider the following code fragment:

1. \(i := 10;\)
2. \(j := 0;\)
3. \(\text{while } (i \geq 0) \{ \}
4. \(i := i - 1;\)
5. \(j := j + 1;\)
6. \(\}\)

and an abstraction \(\rho = \rho_{\text{SIGN}}\), that is, the property of interest is the sign of both \(i\) and \(j\). By computing the abstract semantics of this simple program, we can observe that inside the loop, we lose the sign of \(i\) since \(i\) starts being positive, but then the \(i - 1\) operation makes it impossible to know statically the sign of \(i\) (the result may be positive, zero, or negative starting from \(i\) positive or zero), while we have that \(j\) always remains positive. Moreover, if the loop terminates, we can surely say that, at the end, \(i < 0\); namely, it is negative (due to the negation of the while guard). Hence, we are able to infer that \(i\) is negative and \(j\) is positive after line 6. This means that the final abstract state \(\sigma^{\rho}\) is such that \(\sigma^{\rho}(i) = [\text{neg}]\) and \(\sigma^{\rho}(j) = [\text{pos}]\) (in the following, the extensional notation for \(\sigma^{\rho}\) will be \(\{i \leftarrow [\text{neg}], j \leftarrow [\text{pos}]\}\), similar to the notation for concrete states).

### 3. PROGRAM SLICING

Program slicing [Weiser 1984] is a program manipulation technique that extracts from programs those statements that are relevant to a particular portion of a computation. In order to answer the question about which are the relevant statements, an observer needs a **window** through which only a part of the computation can be seen [Binkley and Gallagher 1996]. Usually, what identifies the portion of interest in the computation is the value of some set of variables at a certain program point, so that a **program slice** comes to be the subset (syntactically, in terms of statements) of the original program that contributes directly or indirectly to the values assumed by some set of variables at the program point of interest. The **slicing criterion** is what specifies the part of the computation that is relevant to the analysis; in this case, a criterion is a pair consisting of a set \(X\) of variables and a program point (or line number) \(n\). The following definition [Binkley and Gallagher 1996] is a possible formalization of the original idea of program slicing [Weiser 1984], in the case of a single variable:

**Definition 3.1.** Binkley and Gallagher [1996]: For a statement \(s\) (at program point \(n\)) and a variable \(x\), the slice \(P'\) of the program \(P\) with respect to the slicing criterion \(\langle s, \{x\} \rangle\) is any executable program with the following properties:

1. \(P'\) can be obtained by deleting zero or more statements from \(P\).
2. If \(P\) halts on the input \(I\), then each time \(s\) is reached in \(P\), it is also reached in \(P'\), and the value of \(x\) at \(s\) is the same in \(P\) and in \(P'\). If \(P\) fails to terminate, then \(s\) may be reached more times in \(P'\) than in \(P\), but \(P\) and \(P'\) have the same value for \(x\) each time \(s\) is executed by \(P\).

It is worth noting that Reps and Yang [1989], in their **slicing theorem**, provide implicitly a similar definition of program slicing, but it only considers terminating computations. The following example provides the intuition of how slicing works.
Consider the word-count program [Majumdar et al. 2007] given in Figure 5. It takes in a block of text and outputs the number of lines \((nl)\), words \((nw)\), and characters \((nc)\). Suppose the slicing criterion only cares for the value of \(nl\) at the end of the program; then a possible slice is on the left in Figure 6. On the other hand, if the criterion is only interested in \(nw\), then a correct slice is on the right.

Starting from the original definition [Weiser 1984], the notion of slicing has gone through several generalizations and versions, but one feature is constantly present: the fact that slicing is based on a notion of semantic equivalence that has to hold between a program and its slices, or on a corresponding notion of dependency, determining what we keep in the slice while preserving the equivalence relation. What we can observe about definitions of slicing such as the one given in Definition 3.1 is that they are precise enough for finding algorithms for soundly computing slicing, such as Reps and Yang [1989], but not formal enough to become suitable to generalizations allowing us to compare different forms of slicing and/or to define new weaker forms of slicing.

In the following, we use the formal framework proposed in Binkley et al. [2006a], where several notions and forms of slicing are modeled and compared. This is not the only attempt to provide a formal framework for slicing (see Section 8), but we believe that, due to its semantic-based approach, it is suitable to include an abstraction level to slicing, which can be easily compared with all the other forms of slicing included in the original framework. Hence, in the following section, we do not rewrite a formal framework; rather, we reformalize the notion of the slicing criterion in order to easily include abstractions as a new parameter. A brief introduction of the formal framework

```java
1 int c, nl := 0, nw := 0, nc := 0;
2 int in := false;
3 while ((c=getchar())!=EOF) {
4   nc := nc+1;
5   if (c=' ' || c='\n' || c='\t') {
6     in := false;
7   } else if (in = false) {
8     in := true;
9     nw := nw+1;
10    } if (c = '\n') {
11      nl := nl+1;
12    }
}
```

Fig. 5. Word-count program.

```java
1 int c, nw := 0;
2 int in := false;
3 while ((c=getchar())!=EOF) {
4   if (c=' ' || c='\n' || c='\t') {
5     in := false;
6   } else if (in = false) {
7     in := true;
8     nw := nw+1;
9       }
10    } if (c = '\n') {
11      nl := nl+1;
12    }
```

Fig. 6. Slices of the word-count program.

**Example 3.2.** Consider the word-count program [Majumdar et al. 2007] given in Figure 5. It takes in a block of text and outputs the number of lines \((nl)\), words \((nw)\), and characters \((nc)\). Suppose the slicing criterion only cares for the value of \(nl\) at the end of the program; then a possible slice is on the left in Figure 6. On the other hand, if the criterion is only interested in \(nw\), then a correct slice is on the right.

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together with some examples showing the differences between the different forms of slicing introduced in the following is given in the appendix.

### 3.1. Defining Program Slicing: The Formal Framework

In this section, our aim is to define the forms of slicing that can be lifted to the abstract level. Namely, we consider an existing framework [Binkley et al. 2006a, 2006b] that allows defining abstract slicing simply by defining an abstract criterion, which, independently from the kind of slicing (static, dynamic, conditional, standard, etc.), observes properties instead of concrete (exact) values. To this end, we need to slightly revise the construction in order to provide a completely unified notation for slicing criteria. The present article will only deal with backward slicing, where the interest is on the part of the program that affects the observation associated with the slicing criterion and not on the part of the program that is affected by such an observation (called instead forward slicing [Tip 1995]).

**Defining Slicing Criteria.** The slicing criterion characterizes what we have to observe of the program in order to decide whether a program is a slice or not of another program. In particular, we have to fix which computations have to be compared, that is, the inputs and the observations on which the slice and the program have to agree.

In the seminal Weiser approach, given a set of variables of interest \( X \) and program statement \( s \), here referred by the program point \( n \) where \( s \) is placed, a slicing criterion was modeled as \( C = (X, n) \). In the following, we will gradually enrich and generalize this model in order to include several different notions and forms of slicing. Weiser’s approach is known as **static slicing** since the equivalence between the original program and the slice has, implicitly, to hold for every possible input. On the other hand, Korel and Laski [1988] proposed a new technique called **dynamic slicing**, which only considers one particular computation, and therefore one particular input, so that the dynamic slice only preserves the (subset of the) meaning of the original program for that input. Hence, in order to characterize a slicing criterion also including dynamic slicing, we have to add a parameter describing the set of initial memories \( I \subseteq M \): the criterion is now \( C = (I, X, n) \), where \( I = M \) for static slicing, while \( I = \{ \mu \} \), with \( \mu \in M \), for dynamic slicing. Finally, Canfora et al. [1998] proposed **conditioned slicing**, which requires that a conditioned slice preserves the meaning of the original program for a set of inputs satisfying one particular condition \( \phi \). Let \( I = \{ \mu \in M \mid \mu \text{ satisfies } \phi \} \) be the set of input memories satisfying \( \phi \) [Binkley et al. 2006a]. Hence, the slicing criterion still can be modeled as \( C = (I, X, n) \).

Each type of slicing comes in four forms, which differ on what the program and the slices must agree on, namely, on the observable semantics that have to agree. In the following, we provide an informal definition of these forms in order to provide the intuition of what will be formally defined afterward:

**Standard:** It considers one point in a program with respect to a set of variables. In other words, the standard form of slicing only tracks one program point. Semantically, this form of slicing consists of comparing the program and the slices in terms of the (denotational) I/O semantics from the program inputs selected by the criterion. Namely, for each selected input, the results of the criterion variables in the point of observation must be the same, independently from the executed statements.

**Korel and Laski** (KL): It is a stronger form where the program and the slice must follow identical paths [Korel and Laski 1988]. Semantically, we could say that the program and the slice must have the same (operational) trace semantics w.r.t. the statements kept in the slice, starting from the program inputs selected by the criterion. In other
words, as before, the final value must be the same, but in this case these values must be obtained by executing precisely the same statements; that is, the sequence of statements involved in the execution of the slice is the same as the original one, apart from the fact that removed statements are missing (see an example at the beginning of the appendix).

Iteration count (IC): When considering the trace semantics, the same program point inside a loop may be visited more than once; in the following, we call the kth iteration of a program point n the kth time the program point n is visited. The iteration count form of slicing requires that a program and its slice agree only at a particular kth iteration of a program point of interest. In this way, when a point of interest is inside a loop, we have the possibility to require that the variables must agree only at some iterations of the loop and not always.

Korel and Laski iteration count (KL): It is the combination of the last two forms.

In order to deal with these different forms of slicing, the slicing criterion must be enriched with additional information. In particular, the KL form of slicing does not change where to observe variables, but it does change the observed semantics up to that point. Hence, we simply have to add a Boolean parameter \( \psi \): true means that we are considering a KL form and we require that the slice must agree with the program on the execution of statements that are in the slice (and obviously also in the original program); on the other hand, false indicates a standard, non-KL form of slicing. Hence, a criterion \( C \) comes to be \( (I, X, n, \psi) \).

The IC form, instead, affects the observation: in order to embed this feature in the criterion, the third parameter has to be changed. Let \( \{k_1, \ldots, k_j\} \subseteq \mathbb{N} \) be the iterations of the program point \( n \in L \) we are interested in; then, instead of \( n \), in the third parameter of the criterion we should have \( \langle n, \{k_1, \ldots, k_j\} \rangle \). Therefore, \( C \) takes the form \( (I, X, O, \psi) \), where \( O \in L \times \wp(\mathbb{N}) \). Note that \( \langle n, \mathbb{N} \rangle \) represents the fact that we are interested in all occurrences of \( n \), as it happens in the standard form.

There are also some simultaneous (SIM) forms of slicing that consider more than one program point of interest. In order to deal with SIM forms of slicing, we simply extend the definition of a slicing criterion by considering \( O \) as a set instead of a singleton, namely, \( O \in \wp(L \times \wp(\mathbb{N})) \).

In the appendix, there are some simple examples showing the main differences between the several forms of slicing introduced so far.

4. ABSTRACT PROGRAM SLICING

This section defines a weaker notion of slicing based on abstract interpretation. In particular, we generalize the existing formal framework discussed in Section 3.1 in order to include abstract versions of slicing.

Program slicing is used for reducing the size of programs to analyze. Nevertheless, this reduction may be insufficient in order to really improve an analysis. Suppose that some variables at some point of execution do not have a desired property (e.g., that they are different from 0, or from null); in order to understand where the error occurred, it would be useful to find those statements that affect such a property of these variables. Standard slicing may return too many statements, making it hard for the programmer to realize which one caused the error.

Example 4.1. Consider the following program \( P \), which inserts a new element elem at position pos in a single-linked list. For simplicity, let pos never exceed the length of list.
Suppose that \texttt{list} is cyclic after line 47; that is, a traversal of the list visits the same node twice. A close inspection of the code reveals that no cycle is created between lines 34 and 47: \texttt{list} is cyclic after line 47 if and only if it was cyclic before line 34.

In the standard approach, it is possible to set the value of \texttt{list} after line 47 as the slicing criterion. In this case, since \texttt{list} can be modified at lines 41 to 47, at least this piece of code must be included in the slice.

On the other hand, let the cyclicity of \texttt{list} after line 47 be the property of interest, represented by \( \rho_{\text{CYC}} \) (Example 2.2). Since this property of \texttt{list} does not change, the entire code can be removed from the slice.

4.1. Defining Abstract Program Slicing

We introduce abstract program slicing, which compares a program and its abstract slices by considering properties instead of exact values of program variables. Such properties are represented as abstract domains, based on the theory of abstract interpretation (Section 2.2).

We first introduce the notion of abstract slicing criterion, where the property of interest is also specified. For the sake of simplicity, the definition only refers to non-SIM forms (i.e., \( \mathcal{O} \) is a singleton instead of a set of occurrences: \( \mathcal{O} \in \mathcal{L} \times \wp(\mathbb{N}) \)). In order to make abstract the criterion, we have to formalize in it the properties we want to observe on program variables. In particular, we could think of observing different properties for different variables. We define a criterion abstraction \( \mathcal{A} \) defined as a tuple of abstract domains, each one relative to a specific subset of program variables: let \( X \) be a set of variables of interest in \( P \) and \( \{X_i\}_{i \in \{1,k\}} \subseteq \wp(\mathcal{X}) \) a partition of \( \mathcal{X} \), and the notation \( \mathcal{A} = (X_1 : \rho_1, \ldots, X_k : \rho_k) \) means that each uco \( \rho_i \) is applied to the set of variable \( X_i \) (left implicit when it is clear from the context), meaning that \( \rho_i \) is precisely the property to observe on \( X_i \). In the following, we denote by \( \mathcal{A}|_{X_i} \) the property observed on \( X_i \), formally \( \mathcal{A}|_{X_i} = \rho_i \). This is the most general representation, accounting also for relational domains. When ucos will be applied to singletons, the notation will be simplified (\( x : \rho \) instead of \( \{x\} : \rho \)).

Example 4.2. Let \( x, y, z, \) and \( w \) be the variables in \( \mathcal{X} \). Let \( \mathcal{A} = (x : \rho_{\text{PAR}}, \{y,z\} : \rho_{\text{INT}}^+, w : \rho_{\text{SIGN}}) \), meaning that the interest is on the parity of \( x \), the sign of \( w \), and the (relational) property of intervals [Cousot and Cousot 1979] of the value \( y + z \). When abstracting a criterion w.r.t. \( \mathcal{A} \), the required observation at a program state \( \sigma \) is

\[
\rho_{\text{PAR}}(\sigma(x)) \quad \rho_{\text{INT}}^+(\sigma(y) + \sigma(z)) \quad \rho_{\text{SIGN}}(\sigma(w)).
\]

In order to be as general as possible, we consider relational properties of variables (see Section 2), so that properties are associated with tuples instead of single variables.
In this case, a property is said to involve some set (tuple) of variables. Given a memory \( \mu \), \( \rho(\mu) \) is the result of applying \( \rho \) to the values in \( \mu \) of the variables involved by the abstract domain, and \( A(\mu) \) is the corresponding notion for tuples of ucos.

**Definition 4.3 (Abstract Criterion).** Let \( I \subseteq M \) be a set of input memories, \( X \subseteq X \) be a set of variables of interest, \( O \in L \times \varphi(N) \) be a set of occurrences of interest, and \( \psi \) be a truth value indicating if the slicing is in KL form. Moreover, let \( X \) be the set of variables of interest and \( A \equiv (X_1 : \rho_1, \ldots, X_k : \rho_k) \), with \( \{X_i\}_{i \in [1,k]} \) a partition of \( X \). Then, the abstract slicing criterion is \( C_A = (I, X, O, \psi, A) \).

Note that, when dealing with nonabstract notions of slicing, we have that each domain is the identity on each single variable, namely, \( \rho = \langle x_1 : \rho_{x_1}, \ldots, x_k : \rho_{x_k} \rangle \), where \( \rho_{x_0} \equiv \lambda x.x \). It is also worth pointing out that, exactly as it happens for nonabstract forms, \( I = M \) corresponds to static slicing, and \( |I| = 1 \) corresponds to dynamic slicing; in the intermediate cases, we have conditioned slicing.

### 4.2. The Extended Formal Framework

In this section, we extend a formal framework in which all forms of abstract slicing can be formally represented. It is an extension of the mathematical structure introduced by Binkley et al. [2006a]. Following their framework, we represent a form of abstract slicing by a pair \((\subseteq, E_A)\), where \( \subseteq \) is the traditional syntactic ordering, and \( E_A \) is a function mapping abstract slicing criteria to semantic equivalence relations on programs. Given two programs \( P \) and \( Q \), and an abstract slicing criterion \( C_A \), we say that \( Q \) is a \((\subseteq, E_A)\)-(abstract)-slice of \( P \) with respect to \( C_A \) if and only if \( Q \subseteq P \) and \( (P, Q) \in E_A(C_A) \) (i.e., \( P \) and \( Q \) are equivalent w.r.t. \( E_A \)). Some preliminary notions are needed to define \( E_A \) in the context of abstract slicing.

An abstract memory w.r.t. a set of variables of interest \( X \) (partitioned in \( \{X_i\}_{i \in [1,k]} \)) is obtained from a memory by restricting its domain to \( X \) and assigning to each set \( X_i \) of variables an abstract value determined by the corresponding abstract property of interest \( \rho_i \).

**Definition 4.4.** Let \( \mu \in M \) be a memory, \( X \) be the set of a tuple of sets of variables of interest, and \( A = (X_1 : \rho_1, \ldots, X_k : \rho_k) \) be the corresponding tuple of properties of interest such that \( \{X_i\}_{i \in [1,k]} \) is a partition of \( X \). The abstract restriction of a memory \( \mu \) w.r.t. the state abstraction \( A \) is defined as \( \mu \lceil_A X \equiv \mu \circ \mu(X) \equiv (\rho_1(\mu(X_1)), \ldots, \rho_k(\mu(X_k))) \).

**Example 4.5.** Let \( X = \{x_1, x_2, x_3, x_4\} \) be a set of variables, and suppose that the properties of interest are the (relational) sign of the product \( x_1 x_2 \) and the parity of \( x_3 \) (both defined in Section 2). We slightly abuse notation by denoting as \( \rho_{\text{SIGN}} \) also its extension to pairs \((u, t)\) where the sign of their product matters: for example, \( \rho_{\text{SIGN}}(3, -5) = (\text{neg}) \). In our formal framework, \( A \) is defined as \( (x_1, x_2) : \rho_{\text{SIGN}}, x_3 : \rho_{\text{PAR}} \). Let \( \mu(x_1) = 1 \), \( \mu(x_2) = 2 \), \( \mu(x_3) = 3 \), and \( \mu(x_4) = 4 \); then, \( \mu \lceil_A X \) comes to be \( \mu \circ \mu(X) = ((\text{pos}), (\text{odd})) \).

The abstract projection operator modifies a state trajectory by removing all those states that do not contain occurrences or points of interest. If there is a state that contains an occurrence of interest, then its memory state is restricted via \( \lceil u \) to the variables of interest, and only a property is considered for each tuple. In the following, the abstract projection \( \text{Proj}^u \) is formally defined.

**Definition 4.6 (Abstract Projection).** Let \( C_A = (I, X, O, \psi, A) \), and \( L \subseteq L \) such that \( L \neq \emptyset \) if \( \psi = \text{true} \), and \( L = \emptyset \) if \( \psi = \text{false} \). For any \( n \in N \), \( k \in N \), \( \mu \in M \), we define a
\[\text{Proj}^\alpha_{(X,O,L,A)}(n^k_1, \mu_1) \cdots (n^k_l, \mu_l)) = \text{Proj}^\alpha_{(X,O,L,A)}(n^k_1, \mu_1) \circ \cdots \circ \text{Proj}^\alpha_{(X,O,L,A)}(n^k_l, \mu_l).\]

The abstract projection \(\text{Proj}^\alpha\) is the extension of \(\text{Proj}^{\alpha_X}\) to sequences:

\[\text{Proj}^\alpha_{(X,O,L,A)}(\langle n^k_1, \mu_1 \rangle, \ldots, \langle n^k_l, \mu_l \rangle) = \text{Proj}^\alpha_{(X,O,L,A)}(n^k_1, \mu_1) \circ \cdots \circ \text{Proj}^\alpha_{(X,O,L,A)}(n^k_l, \mu_l).\]

\text{Proj}^{\alpha_X}\) takes a state from a state trajectory and returns either one pair or an empty sequence \(\varepsilon\). Abstract projection allows us to define all the semantic equivalence relations we need for representing the abstract forms of slicing.

**Example 4.7.** Consider the program \(P\) in Figure 7. Consider \(I = M\) (meaning that we are considering static slicing), \(X = \{i, s\}\), and \(O = (8, N)\) (meaning that we check the value of variables of interest at each iteration of program point 8). Moreover, we consider \(A = \{i : \rho_{\text{SIGN}}, s : \rho_{\text{PAR}}\}\). Then, in Figure 8, we have the corresponding abstract projection (the concrete trace is given in the appendix in Example 24). In this figure, we depict states as tuples of boxes: the first box contains the number of the executed program point (with the iteration counter as apex), while the other boxes indicate the values of each variable. A cross on a box means that the projection does not consider that variable or state. For instance, in this example, we only care about states 6 and 8; moreover, in states 6, we are not interested in the values of variables, while in states 8, we are interested in the sign of \(i\) and in the parity of \(s\) (if we would be interested in the value of these variables we would have the value instead of their property, as it happens in the examples in the appendix).

### 4.3. Abstract Unified Equivalence

The only missing step for completing the formal definition of abstract slicing in the formal framework is to characterize the functions mapping abstract slicing criteria to abstract semantic equivalence relations.
Definition 4.8 (Abstract Unifying Equivalence). Let \( P \) and \( Q \) be executable programs, and \( C_A = (I, X, O, \psi, A) \) be an abstract criterion. Then \( P \) is abstract-equivalent to \( Q \) if and only if, for every \( \mu \in I \), it holds that \( \text{Proj}^{\alpha}_{(X, O, \emptyset, A)}(\tau^\mu_P) = \text{Proj}^{\alpha}_{(X, O, \emptyset, A)}(\tau^\mu_Q) \), where \( L = L_P \cap L_Q \) if \( \psi = \text{true} \). The function \( E_A \) maps each criterion \( C_A \) to a corresponding abstract semantic equivalence relation.

Therefore, a generic form of slicing can be represented as \( (\subseteq, E_A) \). This can be used to formally define both traditional and abstract forms of slicing in the presented abstract formal framework, so that the latter comes to be a generalization of the original formal framework. The following examples show how it is possible to use these definitions in order to check whether a program is an abstract slice of another one.

Example 4.9. Consider programs \( P \) and \( Q \) in Figure 9. Let \( C_A = (M, \{s\}, \{\langle 7, N \rangle \}, \text{false}, \langle s : \rho_{\text{par}} \rangle) \), meaning that we are interested in the parity of \( s \) \( (A = \langle s : \rho_{\text{par}} \rangle) \) at the end of execution \( (O = \{\langle 7, N \rangle \}) \) for all possible inputs \( (I = M) \), in non-KL form. Since \( Q \subseteq P \), in order to show that \( P \) is an abstract static slice of \( Q \) with respect to \( C_A \), we have to show that \( (P, Q) \in E_A(C_A) \) holds. Let \( \mu = \{n \leftarrow a, s \leftarrow b\} \) for some \( a, b \in \mathbb{N} \) be an initial memory. Then, the trajectory of \( P \) from \( \mu \) contains the following computation steps:

![Computation steps](image)

Applying \( \text{Proj}^{\alpha} \) (with \( L = \emptyset \) since \( \psi = \text{false} \)) to \( \tau^\mu_P \) returns only the abstract value of the variable \( s \) at point 7 (due to \( C_A \)):

\[
\text{Proj}^{\alpha}_{(X, O, \emptyset, A)}(\tau^\mu_P) = \langle 7^1, \{n \leftarrow a, s \leftarrow b + a(a + 1), i \leftarrow a + 1\} \mid \alpha \{\langle s \rangle\} \rangle = \langle 7^1, \rho_{\text{par}}(b + a(a + 1)) \rangle = \langle 7^1, \rho_{\text{par}}(b) \rangle.
\]

Since we have \( 7^1 = \langle 7, 1 \rangle \in O \), \( \text{Proj}^{\alpha}_{(X, O, \emptyset, A)}(\tau^\mu_P) = \langle 7^1, \{n \leftarrow a, s \leftarrow b + a(a + 1), i \leftarrow a + 1\} \mid \alpha \{\langle s \rangle\} \rangle \) returns \( 7^1, \mu \mid \alpha X \) = \( 7^1, \{n \leftarrow a, s \leftarrow b + a(a + 1), i \leftarrow a + 1\} \mid \alpha \{\langle s \rangle\} \rangle \). The abstract memory restricts the domain of \( \mu \) to variables of interest, so we consider only the part of \( \mu \) regarding \( s \), that is, \( b + a(a + 1) \). Hence, we have \( 7^1, \rho_{\text{par}}(b + a(a + 1)) \). Since the parity of \( b + a(a + 1) \) only depends on the parity of \( b \), being either \( a \) or \( a + 1 \) even, the final result is \( 7^1, \rho_{\text{par}}(b) \). Consider now the execution of \( Q \) from \( \mu \), which corresponds to the following state trajectory:

![State trajectory](image)

Applying \( \text{Proj}^{\alpha} \) to \( \tau^\mu_Q \) gives:

\[
\text{Proj}^{\alpha}_{(X, O, \emptyset, A)}(\tau^\mu_Q) = \langle 7^1, \{n \leftarrow a, s \leftarrow b\} \mid \alpha \{\langle s \rangle\} \rangle = \langle 7^1, \rho_{\text{par}}(b) \rangle.
\]

ACM Transactions on Computational Logic, Vol. 18, No. 1, Article 7, Publication date: February 2017.
Therefore, \( \text{Proj}_{(X, \mathcal{O}, \emptyset, A)}^\mu(\tau_P^n) \) is equal to \( \text{Proj}_{(X, \mathcal{O}, \emptyset, A)}^\mu(\tau_Q^n) \). As \( \mu \) is an arbitrary input, this equation holds for each \( \mu \in \mathcal{M} \), so that \( \langle P, Q \rangle \in \mathcal{E}(C_A) \), and this implies that \( Q \) is an abstract static slice of \( P \) w.r.t. \( C_A \).

**Example 4.10.** Consider the programs \( R \) and \( S \) in Figure 10, and let \( C_A \) be \((\mathcal{I}, \{s\}, \{(7, \mathbb{N})\}, \text{false}, \{s : \text{par}\})\), where \( \mathcal{I} = \{ \mu \mid \mu(n) \in 4\mathbb{Z} \} \); that is, we are interested in the parity of \( s \) at the end of the execution for all inputs where \( n \) is a multiple of 4. Since \( S \subseteq R \), in order to show that \( S \) is an abstract conditioned slice of \( R \) w.r.t. \( C_A \), we have to show that \( \langle R, S \rangle \in \mathcal{E}(C_A) \) holds, namely, that they have the same abstract projection. Let \( \mu \in \mathcal{I} \) be an initial memory, and suppose \( \mu(n) = a = 4m \). The trajectory \( \tau_R^\mu \) of \( R \) from \( \mu \) contains the following steps:

\[
\begin{array}{cccccccc}
\text{a} & 1 & 2 & 3 & 4 & 5 & 6 & 7 & 8 \\
\text{n} & s & a & b & a & b & a+1 & a & a+1 \\
\end{array}
\]

On the other hand, executing \( S \) from \( \mu \) gives the state trajectory \( \tau_S^\mu \).

Applying \( \text{Proj}^\mu \) to both state trajectories, we have:

\[
\text{Proj}_{(X, \mathcal{O}, \emptyset, A)}^\mu(\tau_R^\mu) = \langle 7^1, \{n \leftarrow a, s \leftarrow a(a+1)/2, i \leftarrow a+1 \} \mid_A \langle \{s\} \rangle \rangle = \langle 7^1, \rho_{\text{par}}(a(a+1)/2) \rangle
\]

\[
= \langle 7^1, \rho_{\text{par}}(2m(4m+1)) \rangle = \langle 7^1, 2\mathbb{Z} \rangle
\]

\[
= \langle 7^1, \rho_{\text{par}}(0) \rangle = \langle 7^1, \{n \leftarrow 0, s \leftarrow 0 \} \mid_A \langle \{s\} \rangle \rangle = \text{Proj}_{(X, \mathcal{O}, \emptyset, A)}^\mu(\tau_S^\mu).
\]

Therefore, we have \( \text{Proj}_{(X, \mathcal{O}, \emptyset, A)}^\mu(\tau_R^\mu) = \text{Proj}_{(X, \mathcal{O}, \emptyset, A)}^\mu(\tau_S^\mu) \). Since \( \mu \) is an arbitrary input from \( \mathcal{I} \), this equation holds for each \( \mu \in \mathcal{I} \), so that \( \langle R, S \rangle \in \mathcal{E}(C_A) \), and this implies that \( S \) is an abstract conditioned slice of \( R \) with respect to \( C_A \). It is worth noting that \( S \) is not a static abstract slice of \( R \) since, for all the input values \( a \notin 4\mathbb{Z} \) for \( n \), the parity of the final value of \( s \) is not necessarily even.

### 4.4. Comparing Forms of Abstract Slicing

This section provides a formal theory allowing us to compare abstract forms of slicing between themselves and with nonabstract ones. First of all, we show under which conditions an abstract semantic equivalence relation *subsumes* another one; analogously, we show when the form of (abstract) slicing, corresponding to the former equivalence relation, subsumes the form of (abstract) slicing corresponding to the latter one. Such results are necessary in order to obtain a precise characterization of the extension of the weaker than relation (whose original definition is recalled in the appendix) to the abstract forms of slicing.

The following lemma shows under which conditions on the slicing criteria there is a relation of subsumption between two semantic equivalence relations. In the
following, we denote with \( \sqsubseteq \) the relation “more concrete than” in the lattice of abstract interpretations between tuples of abstractions. Formally, let us consider \( \mathcal{A}^1 = (X_1^1: \rho_1^1, \ldots, X_n^1: \rho_n^1) \) defined on the variables \( \chi^1 \) and \( \mathcal{A}^2 = (X_1^2: \rho_1^2, \ldots, X_n^2: \rho_n^2) \) defined on the variables \( \chi^2 \), such that \( \chi^1 \subseteq \chi^2, k_1 \leq k_2 \), and \( \forall i \leq k_1 \); we have \( X_i^1 = X_i^2 \), namely, that the variables in common are partitioned in the same way. Then \( \mathcal{A}^2 \sqsubseteq \mathcal{A}^1 \) if and only if \( X_i^1 \in \chi^1, \rho_i^2 \subseteq \rho_i^1 \). Note that, for all the variables in \( \chi^2 \setminus \chi^1 \), the abstraction \( \mathcal{A}^1 \) does not require any particular observation; hence, on these variables, surely \( \mathcal{A}^2 \) is more precise. The following relation is such that, when both \( \mathcal{A}^1 \) and \( \mathcal{A}^2 \) are the identity on all the variables of interest, then the resulting criterion relation is the same as proposed in Binkley et al. [2006b] (see the appendix for details) for characterizing the original formal framework.

**Lemma 4.11.** Let two abstract slicing criteria \( \mathcal{C}_A^1 = (\mathcal{I}^1, \chi^1, \mathcal{O}^1, \psi^1, \mathcal{A}^1) \) and \( \mathcal{C}_A^2 = (\mathcal{I}^2, \chi^2, \mathcal{O}^2, \psi^2, \mathcal{A}^2) \) be given. If (1) \( \mathcal{I}^1 \subseteq \mathcal{I}^2 \), (2) \( \mathcal{O}^1 \sqsubseteq \mathcal{O}^2 \), (3) \( \chi^1 \subseteq \chi^2 \), (4) \( \psi^1 \Rightarrow \psi^2 \), and (5) \( \mathcal{A}^2 \sqsubseteq \mathcal{A}^1 \) (denoted \( \mathcal{C}_A^1 \rightarrow \mathcal{A} \mathcal{C}_A^2 \)), then \( (\sqsubseteq, \mathcal{E}(\mathcal{C}_A^1)) \) subsumes \( (\sqsubseteq, \mathcal{E}(\mathcal{C}_A^2)) \); that is, for every \( P \) and \( Q \) such that \( Q \sqsubseteq P \), \( (P, Q) \in \mathcal{E}(\mathcal{C}_A^2) \) implies \( (P, Q) \in \mathcal{E}(\mathcal{C}_A^1) \).

**Proof.** First of all, note that, if \( \mathcal{A}^1 = \mathcal{A}^2 = \rho_{id} \), namely, if we are considering concrete criteria, then \( \rightarrow \mathcal{A} \) collapses to the concrete relation defined in Binkley et al. [2006a] (which is the \( \rightarrow \) defined in Equation (2) in the Appendix). Hence, in this case, the results hold by Binkley et al. [2006a]. Suppose \( (P, Q) \in \mathcal{E}(\mathcal{C}_A^2) \) with \( Q \sqsubseteq P \), namely, \( Q \) slice of \( P \) w.r.t. \( \mathcal{E}(\mathcal{C}_A^2) \). This means that, for each \( \mu_0 \in \mathcal{I}^2 \), \( \mathcal{P}_Q^{\mathcal{O}_2, \mathcal{C}_A^2}(\tau_Q^{\mathcal{O}_0}) = \mathcal{P}_Q^{\mathcal{O}_2, \mathcal{C}_A^2}(\tau_Q^{\mathcal{O}_0}), \) where \( \mathcal{L}^2 \) is defined as in Definition 4.8. This means that, for each state \( (n^k, \mu) \) in the trajectory \( \tau_Q^{\mathcal{O}_0} \), whose projection \( \mathcal{P}_Q^{\mathcal{O}_2, \mathcal{C}_A^2}(n^k, \mu) \) is not empty, there exists a state \( (n^k, \mu) \) in \( \tau_Q^{\mathcal{O}_0} \) with the same projection. Let us consider now, \( \mu \in \mathcal{I}^1 \subseteq \mathcal{I}^2 \). We prove that on these states \( \mathcal{P}_Q^{\mathcal{O}_2, \mathcal{C}_A^2}(n^k, \mu) = \mathcal{P}_Q^{\mathcal{O}_2, \mathcal{C}_A^2}(n^k, \mu) \) (and in this case there is a corresponding state in the trajectory of \( Q \), or it is empty (and in this case also the state in \( Q \) has empty projection). Recall that

\[
\mathcal{P}_Q^{\mathcal{O}_2, \mathcal{C}_A^2}(n^k, \mu) \triangleq \begin{cases} 
(n^k, \mu |_{\mathcal{A}^1} \chi^1) & \text{if } \exists (n, K) \in \mathcal{O}^1 \text{. } k \in K \\
(n^k, \bot) & \text{if } \nexists (n, K) \in \mathcal{O}^1 \text{. } k \in K \text{ and } n \in \mathcal{L}^1 
\end{cases}
\]

where also \( \mathcal{L}^1 \) is defined as in Definition 4.8. Note that, since we are considering both the criteria on the same pair of programs, we have also that \( \psi^1 \Rightarrow \psi^2 \) corresponds to saying that \( \mathcal{L}^1 \subseteq \mathcal{L}^2 \). At this point:

- If \( \mathcal{P}_Q^{\mathcal{O}_2, \mathcal{C}_A^2}(n^k, \mu)_{P} = (n^k, \mu |_{\mathcal{A}^1} \chi^1) \), then \( \exists (n, K) \in \mathcal{O}^1 \text{. } k \in K \), but \( \mathcal{O}^1 \subseteq \mathcal{O}^2 \), and hence \( (n, K) \in \mathcal{O}^2 \text{. } k \in K \). This means that \( \mathcal{P}_Q^{\mathcal{O}_2, \mathcal{C}_A^2}(n^k, \mu)_{P} = (n^k, \mu |_{\mathcal{A}^2} \chi^2) \), which by hypothesis is equal to \( \mathcal{P}_Q^{\mathcal{O}_2, \mathcal{C}_A^2}(n^k, \mu)_{Q} \), for a memory \( \mu. \) By definition and hypothesis, \( \mu |_{\mathcal{A}^2} \chi^2 = \mathcal{A}^2 \circ \mu(\chi) = \mathcal{A}^2 \circ \mu(\chi) \). Namely, \( \rho_1^2(\mu(X_1)), \ldots, \rho_{k^2}(\mu(X_{k^2})) = \rho_1^2(\mu(X_1)), \ldots, \rho_{k^2}(\mu(X_{k^2})) \). Therefore, in particular, \( \forall i \in \{1, k^1\} \subseteq \{1, k^2\} \) and we have \( \rho_i^2(\mu(X_i)) = \rho_i^2(\mu'(X_i)) \), but by hypothesis \( \rho_i^2 \subseteq \rho_i^1 \), and hence we also have \( \rho_i^1(\mu(X_i)) = \rho_i^1(\mu'(X_i)) \) (by properties of ucos). But then \( \rho_1^1(\mu(X_1)), \ldots, \rho_{k^1}(\mu(X_{k^1})) = \rho_1^1(\mu'(X_1)), \ldots, \rho_{k^1}(\mu'(X_{k^1})) \), namely,

\[The notation \( \mathcal{P}_Q^{\mathcal{O}_2, \mathcal{C}_A^2}(n^k, \mu)_{P} \) means that we are projecting a state of the computation of \( P \).]
Let $D$ be semantic equivalence, $\mu$ denotes dynamic hierarchy obtained by enriching the hierarchy in Figure 25 with standard forms of slicing simply by using the comparison notions defined earlier. Nonabstract equivalence relations such that $\rho_\mu$ subsumes $\rho$ and a slice given a program and a slicing criterion. Again, among all the equivalence relations introduced in Section 4.3 subsume the corresponding abstract equivalence relations $\mu \subseteq \rho$. Further, by using this lemma, we can show that $AD(C_A)$ subsumes $AC(C_A)$, which in turns subsumes $AS(C_A)$.

Theorem 4.12 (Binkley et al. [2006b]). Let $R_1$ and $R_2$ be semantic equivalence relations such that $R_2$ subsumes $R_1$. Then, for every $P$ and $Q$, we have $(P, Q) \in (\subseteq, R_1) \Rightarrow (P, Q) \in (\subseteq, R_2)$.

Figure 11 shows the non-SIM hierarchy obtained by enriching the hierarchy in Figure 25 with standard forms of abstract static slicing, abstract dynamic slicing, and abstract conditioned slicing. In general, we can enrich this hierarchy with any abstract form of slicing simply by using the comparison notions defined earlier. Nonabstract forms are particular cases of abstract forms of slicing, as they can be instantiated by choosing the identity property, $\rho_{id}$, for each variable of interest. Hence, nonabstract forms are the “strongest” forms, since, for each property $\rho$, we have $\rho_{id} \subseteq \rho$. Moreover, if parameters $M, X, \mathcal{O}, \psi$ are fixed, and $A$ is made less precise or more abstract (i.e., the information represented by the property is reduced), then the abstract slicing form becomes weaker, as suggested by dotted lines in Figure 11.

5. PROGRAM SLICING AND DEPENDENCIES

The previous sections introduced a formal framework containing different notions of program slicing. In particular, we observed that a kind of slicing is a pair consisting of a syntactic preorder and a semantic equivalence relation [Binkley et al. 2006a]. After discussing how the notion of “to be a slice of” can be formally defined, the focus will shift to how to compute a slice given a program and a slicing criterion. Again, among all the
possible definitions of slicing, we are interested in slices obtained by erasing statements from the original program; that is, the slice is related to the original program by the syntactic ordering relation $\sqsubseteq$. Given a slicing criterion, the idea is to keep all the statements affecting the semantic equivalence relation defined by the criterion. In other words, we should have to translate the formal definition into a characterization of which statements have to be kept in a slice, or vice versa which statements can be erased, in order to preserve the semantic equivalence defining the chosen notion of slicing. Intuitively, we have to keep all the statements affecting the semantics defined by the chosen slicing criterion.

The standard approach for characterizing slices and the corresponding relation being slice of is based on the notion of program dependency graph [Horwitz et al. 1989; Reps and Yang 1989], as described by Binkley and Gallagher [1996]. Program dependency graphs (PDGs) can be built out of programs and describe how data propagate at runtime. In program slicing, we could be interested in computing dependencies on statements: $s''$ depends on $s'$ if some variables that are used inside $s''$ are defined inside $s'$, and definitions in $s'$ reach $s''$ through at least one possible execution path. Also, $s$ depends implicitly on an if-statement or a loop if its execution depends on the Boolean guard.

**Example 5.1.** Consider the program in Figure 12 and the derived PDG (edges that can be obtained by transitivity are omitted). $s_8$ depends on both $s_5$ and $s_7$ (and, by transitivity, $s_1$) since $v$ is not known statically when entering $s_8$. On the other hand, there is no dependency of $s_8$ on either (i) $s_6$, since $z$ is not used in $s_8$, or (2) $s_2$, since $w$
is always redefined before $s_8$. The dependency of $s_7$ on $s_1$ is implicit since 4 does not depend on $x$ or $y$, but $s_7$ is executed conditionally on $s_1$.

Formally, a program dependency graph [Gallagher and Lyle 1991] $G_P$ for a program $P$ is a directed graph with nodes denoting program components and edges denoting dependencies between components. The nodes of $G_P$ represent the assignment statements and control predicates in $P$. In addition, nodes include a distinguished node called $Entry$, denoting where the execution starts. An edge represents either a control dependency or a flow (data) dependency. Control dependency edges $u \rightarrow_c v$ are such that (1) $u$ is the entry node and $v$ represents a component of $P$ that is not nested within any control predicate, or (2) $u$ represents a control predicate and $v$ represents a component of $P$ immediately nested within the control predicate represented by $u$. Flow dependency edges $u \rightarrow_f v$ are such that (1) $u$ is a node that defines the variable $x$ (usually an assignment), (2) $v$ is a node that uses $x$, and (3) control can reach $v$ from $u$ via an execution path along which there is no intervening redefinition of $x$.

Unfortunately, there is clearly a gap between the definition of slicing given in Definition 3.1 and the standard implementation based on PDGs [Horwitz et al. 1989; Reps 1991]. This happens because slicing and dependencies are usually defined at different levels of approximation. In particular, we can note that the slicing definition in the formal framework defines slicing by requiring the same behavior, with respect to a criterion, between the program and the slice; that is, we are specifying what is relevant as a semantic requirement. On the other hand, dependency-based approaches consider a notion of dependency between statements that corresponds to the syntactic presence of a variable in the definition of another variable. In other words, slices are usually defined at the semantic level, while dependencies are defined at the syntactic level. The idea presented in this article consists, first of all, of identifying a notion of semantic dependency corresponding to the slicing definition given earlier, in order to characterize the implicit parametricity of the notion of slicing on a corresponding notion of dependency. This way, we can precisely identify the semantic definition of slicing corresponding to a given dependency-based algorithm, characterizing so far the loss of precision of a given algorithm w.r.t. the semantic definition. In Figure 13, we show these relations. In particular, starting from the criterion, we can define an equivalent notion of dependency that allows us to identify which variables should be kept in a slice, affecting the whole program semantics (program dependency notion).

Definition 5.2 ((Semantic) Program Dependency). Let $C = (\mathcal{I}, \mathcal{X}, \mathcal{O}, \psi)$ be a slicing criterion and $P$ be a program. The program depends on $x$, denoted $x \xrightarrow{\text{C}} P$, if and only if

$$
\exists \mu_1, \mu_2 \in \mathcal{I}. \forall y \neq x. \mu_1(y) = \mu_2(y) \land \text{Proj}_{(\mathcal{X}, \mathcal{O}, \mathcal{L})}(\tau^1_P) \neq \text{Proj}_{(\mathcal{X}, \mathcal{O}, \mathcal{L})}(\tau^2_P).
$$

This means that the variable $x$ affects the observable semantics of $P$.

Unfortunately, this characterization is not effective due to undecidability of the program semantics. In particular, the amount of traces to compare could be infinite, and
also the traces themselves could be infinite. Hence, we consider a stronger notion of dependency that looks for local semantic dependencies, identifying all the variables affecting at least one expression used in the program (expression dependency notion), and this is precisely the semantic generalization of the syntactic dependency notion used, for instance, in a PDG-based algorithm for slicing. In other words, we characterize when a variable affects the semantics of an expression in $P$.

Our idea is to make semantic the standard notion of syntactic dependency, by substituting the notion of uses with the notion of depends on [Giacobazzi et al. 2012]. In order to obtain this characterization, we have to find which variables might affect the evaluation of the expression $e$ in the assignment $z := e$ or in a control statement guarded by $e$, that is, which variables belong to the set $\text{REL}(e)$ of the variables relevant to the evaluation of $e$. As already pointed out, standard syntactic dependency calculi compute $\text{REL}(e)$ as $\text{VARS}(e)$.

**Definition 5.3 ((Semantic) Expression Dependency).** Let $C = (I, X', O, \psi)$ be a slicing criterion. Let $x \in X$, $Y \subseteq X$.

$$
x \xRightarrow{C} e \iff \exists \sigma_1, \sigma_2 \in \Sigma. \forall y \neq x. \sigma_1(y) = \sigma_2(y) \land [e](\sigma_1) \neq [e](\sigma_2)
$$

$$
Y \xRightarrow{C} e \iff \exists y \in Y. y \xRightarrow{C} e
$$

The formulation of $x \xRightarrow{C} e$ can be rewritten as

$$
\exists \sigma \in \Sigma, v_1, v_2 \in \mathbb{V}. [e](\sigma [x \leftarrow v_1]) \neq [e](\sigma [x \leftarrow v_2]).
$$

**Proposition 5.4.** Let $C$ be a slicing criterion. If $x \xRightarrow{C} P$, then there exists $e$ in $P$ such that $x \xRightarrow{C} e$.

**Proof.** By contradiction. If for each $e$ in $P$ we have $x \not\xRightarrow{C} e$, then no expression in $P$ depends on $x$. Therefore, independently from $x$, $P$ provides precisely the same results. □

By using this notion of dependency, we can characterize the subset $\text{REL}(e) \subseteq \text{VARS}(e)$ containing exactly those variables that are semantically relevant for the evaluation of $e$. This way, we obtain a notion of dependency that allows us to derive more precise slices, that is, to remove statements that a merely syntactic analysis would leave.

**Example 5.5.** Consider the program $P$:

```plaintext
1     x := e_x;
2     y := e_y;
3     w := e_w;
4     z := w + y + 2(x^2) - w;
```

where $e_x, e_y$, and $e_w$ are expressions. We want to compute the static slice $Q$ of $P$ affecting the final value of $z$ (i.e., the slicing criterion $C = (M, \{z\}, (4, N), \text{false})$ is interested in the final value of $z$). If we consider the traditional notion of slicing, then it is clear that we can erase line 3 without changing the final result for $z$. However, in the usual syntactic approach, we would have a dependency between $z$ and $w$, since $w$ is used where $z$ is defined. Consequently, the slice obtained by applying this form of dependency would leave the program unchanged. On the other hand, if the semantic dependency is considered, then the evaluation of $w + y + 2(x^2) - w$ does not depend on the possible variations of $w$, which implies that we are able to erase line 3 from the slice.

Next we show how the PDG-based approach to slicing can be modified in order to cope with semantic slicing. The PDG approach is based on the computation of the set of variables used in an expression $e$. In the following, we explore if this set can be rewritten by considering a semantic form of dependency.
Hence, let us define the new notion of semantic PDG, where all the flow dependencies are semantic; that is, we substitute the flow edges defined earlier with semantic flow dependency edges \( u \rightarrow_{sf} v \) that are such that (1) \( u \) is a node that defines the variable \( x \) (usually an assignment), (2) \( v \) is a node containing an expression \( e \) such that \( x \sim_{C} e \) (where \( C \) is the criterion with respect to which we are computing the slice), and (3) control can reach \( v \) from \( u \) via an execution path along which there is no intervening redefinition of \( x \). A (semantic) flow path is a sequence of (semantic) flow edges.

**Proposition 5.6.** Let \( P \) be a program and \( C \) be a slicing criterion. Let \( G_{P} \) be the PDG with flow dependency edges \( \rightarrow_{f} \), and \( G_{P}^{s} \) be the semantic PDG where the flow dependency edges are semantic \( \rightarrow_{sf} \). If \( u \rightarrow_{sf} v \), then \( u \rightarrow_{f} v \).

**Proof.** Trivially, since \( x \sim_{C} e \), then \( e \) must use the variable \( x \). \( \square \)

In principle, a (backward) slice is composed by all the statements (i.e., nodes) such that there exists a path from the corresponding node to the relevant (according to the slicing criterion) use of a variable of interest (in the criterion) [Reps and Yang 1989]. In other words, we follow backward the (semantic) flow edges from the nodes identified by the criterion, and we keep all the nodes/statements we reach. Hence, the criterion defines the edges that we can follow in order to compute the slice.

By using the semantic flow dependency edges, we can draw a new semantic PDG containing fewer flow edges, that is, only those corresponding to semantic dependencies. At this point, the type of slicing (either static, dynamic, or conditional) characterized by the criterion decides which nodes can be kept in the PDG.

**Theorem 5.7.** Let \( \mathcal{C} = (\mathcal{I}, \mathcal{X}, \mathcal{O}, \psi) \) be a slicing criterion. Let \( P \) be a program and \( G_{P}^{s} \) be its semantic PDG, that is, a PDG whose flow edges are \( \rightarrow_{sf} \). Let \( Q \) be the subprogram of \( P \) containing all the statements corresponding to nodes such that there exists a semantic flow path in \( G_{P}^{s} \) from them to a node in \( \mathcal{N}_{\mathcal{O}} = \{ n \mid \exists (n, K) \in \mathcal{O} \} \). Then \( Q \) is a slice w.r.t. the criterion \( \mathcal{C} \).

**Proof.** Note that, in the PDG construction, all the flow edges \( \rightarrow_{sf} \) satisfy a notion of semantic dependencies holding for all inputs by definition, while all the other edges are syntactic, and hence its construction is independent from the input set \( \mathcal{I} \), meaning that any slice computed by using the semantic PDG holds for any possible input memory in \( \mathbb{M} \). Moreover, since we simply collect the statements potentially affecting the program observation, we cannot decide which iteration to observe, and for this reason each program point is taken in the slice independently from the iteration to observe, and for this reason the obtained slice will provide the same result for any possible iteration; that is, the set of interesting points to observe are \( \mathcal{N}_{\mathcal{O}} \times \{ \mathbb{N} \} = \{ (n, \mathbb{N}) \mid n \in \mathcal{N}_{\mathcal{O}} \} \). Finally, by construction, we cannot have statements executed in the slice that are not executed in the original programs. Hence, the criterion enforced by the PDG slice construction is \( C_{PDG} = (\mathbb{M}, \mathcal{X}, \mathcal{N}_{\mathcal{O}} \times \{ \mathbb{N} \}, \text{true}) \) (direct consequence of the slicing theorem in Reps and Yang [1989]), which for each \( \mathcal{I} \subseteq \mathbb{M} \) and \( \{ k_{n} \}_{n \in \mathcal{N}_{\mathcal{O}}} \subseteq \wp(\mathbb{N}) \) is a slice also w.r.t. the criterion \( (\mathcal{I}, \mathcal{X}, \{ (n, K_{n}) \mid n \in \mathcal{N}_{\mathcal{O}} \}, \psi) \), by Binkley et al. [2006a] (Equation (2) in the appendix). Hence, we have that the result is a slice w.r.t. \( C \) [Binkley et al. 2006a]. \( \square \)

We have pointed out so far the difference between syntactic and semantic dependencies: it can be the case that a variable syntactically appears in an expression without affecting its value. Actually, one could argue that the case is not very likely to happen: the possibility to find an assignment like \( x := y \rightarrow y \) in code written by a professional software engineer is remote at the very least. However, when it comes to abstract dependencies, the picture is quite different, and we could even say that, in the present work, (concrete) semantic dependencies have been mainly introduced to prepare the discussion about
their abstract counterpart. Indeed, it is much more likely that some variables are not semantically relevant to an expression if the value of interest is an abstract one, for example, the parity or the sign of a numeric expression, or the nullity of a pointer. This justifies the definition of a semantic notion of dependency at the abstract level.

6. ABSTRACT DEPENDENCIES

This section discusses the problem of defining and computing abstract dependencies, which allow capturing the dependency relation between variables with respect to a given abstract criterion. The previous section formalized this relation in the concrete semantic case; the following example takes it to the abstract level.

Example 6.1. Consider the expression \( e = 2x^2 + y \): although both variables are semantically relevant to the result, only \( y \) can affect its parity, since \( 2x^2 \) will always be even. On the other hand, note that both variables are relevant to the sign of \( e \), in spite of the positivity of \( x^2 \). In fact, given a negative value for \( y \), a change in the value of \( x \) can alter the sign of the entire expression.

The notion of **semantic program dependency** can be easily extended to abstract criteria by changing the projection considered. In this case, we write \( x \xrightarrow{\rho} C_A \), meaning that \( x \) has an effect on the abstract projection of \( P \) determined by \( C_A \). Computing abstract dependencies is still undecidable; hence, again, we need an approximation of the semantic notion.

6.1. Abstract Slicing and Dependencies

Previously, we defined concrete semantic dependencies by identifying those variables that affect the final observation of an expression. Analogously, in order to define abstract semantic dependencies, we need to consider the abstract interference between a property of a variable and a property of an expression.

The next definition below follows the same philosophy as narrow abstract noninterference [Giacobazzi and Mastroeni 2004a; Mastroeni 2013], where abstractions for observing input and output are considered, but these abstractions are observations of the concrete executions.

**Definition 6.2 (Abstract Semantic Expression Dependency, Ndep).** Consider the \( \text{ucos} \rho \in \text{ucos}(\varphi(V)) \) and \( \tilde{\eta} \in \text{ucos}(\varphi(V))^n \), where \( n \) is the number of variables; that is, \( \tilde{\eta} \) is a tuple such that \( \eta_y \) is the property on the variable \( y \). Then:

\[
x \xrightarrow{\text{Ndep}} e \iff \exists \sigma_1, \sigma_2 \in \Sigma . \ (\forall y \neq x . \eta_y(\sigma_1(y)) = \eta_y(\sigma_2(y)) \land \rho[\llbracket e \rrbracket(\sigma_1)) \neq \rho[\llbracket e \rrbracket(\sigma_2))].
\]

This notion is a generalization of Definition 5.3 where we abstract the observation of the result by applying \( \rho \) and the information that we fix about all the variables different from \( x \) by using \( \tilde{\eta} \). Still, this notion characterizes whether the variation of the value of \( x \) affects the abstract evaluation in \( \rho \) of \( e \).

As an important result, we have \( x \xrightarrow{\text{Ndep}} e \iff x \xrightarrow{\text{Ndep}} e \) whenever \( \rho^1 \) and \( \rho^2 \) induce the same partitions (Section 2.3), either on values or tuples of values. This happens because, in Definition 6.2, both abstractions are only applied to singletons. In the following, only partitioning domains will be considered since it is straightforward to note that \( x \xrightarrow{\text{Ndep}} e \) is affected only by \( \Pi(\rho) \), rather than by \( \rho \) itself.

When dealing with abstractions and abstract computations, some more issues have to be taken into account. Consider the program in Example 5.5, and consider the \( \rho_{\text{VAR}} \) property (Section 2.2) for all variables on both input and output. If we compute the set of

---

5For simplicity, we do not consider relational abstract domains in this section.
variables on which the parity of \( e = w + y + 2(x^2) - w \) depends, then we can observe that \( e \) is still independent from \( w \) but is also independent from any possible variation of \( x \). At first sight, the parity of \( w + y + 2(x^2) - w \) is independent from \( x \) just because \( 2(x^2) \) is constantly even. However, it is not only a matter of constancy: a deeper analysis would note that we can look simply at the abstract value of \( x \) only because the operation involved (the sum) in the evaluation is \textit{complete} (Section 2.2), that is, precise, with respect to the abstract domain considered (\( \rho_{\text{par}} \)). In particular, when we deal with abstract domains that are complete for the considered operations, then it is enough to look at the abstract value of variables in order to compute dependencies. Indeed, consider the \( \rho_{\text{sign}} \) domain (Section 2.2): even if the sign of \( 2(x^2) \) is constantly positive, the final sign of \( z \) might be affected by a concrete variation of \( x \) (e.g., consider \( y = -4 \) and two executions in which \( x \) is, respectively, 1 and 5). Therefore, \( x \) has to be considered relevant, although the sign of \( 2(x^2) \) (the only subexpression containing \( x \)) is constant. This can also be derived by considering the logic of independencies from Amtoft and Banerjee [2007] since, by varying the value of \( x \), we can change the sign of \( e \).

Unfortunately, the notion of abstract dependencies given in Definition 6.2 is not suitable for weakening the PDG approach, as we have done in the concrete semantic case. Let us explain the problem in the following example.

\textbf{Example 6.3.} Consider the program \( P = C; x := (y > 0 ? 0 : 1) \), where \( C \) is some code fragment and the expression \( b ? e1 : e2 \) evaluates to \( e1 \) if \( b \) is true, and to \( e2 \) otherwise. Suppose the criterion requires the observation of the parity of \( x \) at this program point, that is, \( A = \langle x : \rho_{\text{par}} \rangle \). Then, it is straightforward to observe that the expression depends on \( y \), but we would like to be more specific (being in the context of abstract slicing), and we can observe that it is the sign of \( y \) that affects the parity of the expression, and therefore of \( x \). Therefore, in the code \( c \), we should look for the variables affecting not simply \( y \) (as expected in standard slicing approaches), but more specifically the \textit{sign} of \( y \), a requirement that is not considered in the abstract criterion.

This example shows that, if we aim at computing abstract dependencies precisely, we would need an algorithm able to keep trace (backward) not only of the different variables that become of interest (i.e., affecting the desired criterion) but also of the different properties to observe on such variables. This means that, while for concrete semantic program dependencies we can provide a definition depending only on the criterion, this is not possible in an abstract context, because each flow edge should be defined depending on abstract properties potentially different from those in the abstract criterion, and which should be characterized dynamically backward starting from the criterion. Unfortunately, this is not possible in Definition 6.2, where we always look for the variation of the value rather than an abstract property of \( x \). Hence, using this notion for substituting the semantic dependency in PDGs would not give the desired results.

These observations make clear that, if we aim at constructively characterizing abstract slicing by means of the abstract dependency notion provided in Definition 6.2, we need to build from scratch a systematic approach for characterizing abstract slicing. The first step in this direction is a computable approximation of the abstract dependencies of Definition 6.2.

\textbf{6.2. A Constructive Approach to Abstract Dependencies}

By means of the (uco-dependent) definition of operations on abstract values, it is possible to automatically obtain (an overapproximation of) the set of relevant variables. The starting point is the \textit{brute-force} approach, which uses the abstract version of concrete operations and explicitly \textit{goes into} the quantifiers involved in Definition 6.2.
Example 6.4. Consider the program in Example 5.5 and let $\eta_y = \rho_{\text{PAR}}$. In order to decide whether $x \rightsquigarrow e$ holds for some $\rho$, the brute-force approach considers the abstract evaluation of $e$ in all contexts where all variables different from $x$ do not change, up to $\overline{\eta}$, while $x$ may change. In this example, this boils down to considering pairs of memories where $y$ has the same parity (with no information about sign), whereas $x$ changes freely, and seeing whether the final values of $e$ agree on $\rho$. Suppose $\rho_{\text{PAR}}(y) = \{\text{even}\}$ (meaning that it is even but the sign is unknown) and $\rho = \rho_{\text{SIG}}$; then, we should compute the abstract value of $e$ for each possible value of $x$. We can easily find $\sigma_1$ and $\sigma_2$ such that

$$\rho_{\text{SIG}}(2 \times \sigma_1(x)^2 + \sigma_1(y)) = \text{[neg]}, \quad \rho_{\text{SIG}}(2 \times \sigma_2(x)^2 + \sigma_2(y)) = \text{[pos]},$$

even if $\rho_{\text{PAR}}(\sigma_1(y)) = \rho_{\text{PAR}}(\sigma_2(y))$ (for instance, $\sigma_1(y) = \sigma_2(y) = -4$ while $\sigma_1(x) = 1$ and $\sigma_2(x) = 5$). It is clear that the sign of $e$ may depend on the value of $x$ since to “fix” the abstract property of the other variables is not enough to “fix” the final value of $e$ with respect to $\rho$. On the other hand, the parity of $e$ does not depend on $x$: if we fix the property $\rho_{\text{PAR}}$ of $y$, for instance, to $\{\text{even}\}$ (but it holds also for the other abstract values), then

$$2 \times \text{[negodd]}^2 + \text{[even]} = \text{[even]}, \quad 2 \times \text{[poseven]}^2 + \text{[even]} = \text{[even]} \quad \ldots.$$

That is to say, if we fix all the (abstract values of the) variables but $x$, then the parity of the result does not change, so that the variation of $x$ does not affect the parity of $e$.

In the following, we introduce an algorithm to improve the computational complexity of the brute-force approach, especially on bigger ucos, and when (1) several variables are involved in expressions and (2) a significant part of them is irrelevant.

6.2.1. Checking Ndep. We discuss how to constructively compute dependencies. Unfortunately, in static analysis, the concrete semantics cannot be used directly as it appears in Definition 6.2, so we need to design an overapproximation. The following definition introduces a stronger notion of dependency based on a sound abstract semantics $\llbracket \cdot \rrbracket^\rho$ (Section 2.4), which approximates $Edep$.

Definition 6.5 (Atom-Dep). An expression $e$ atom-depends on $x$ (written $x \rightsquigarrow e$) with respect to $\rho \in uco(\wp(\Sigma))$ and $\overline{\eta} \in uco(\wp(\Sigma))^n$ ($n$ is the number of variables) if and only if there exist $\sigma_1, \sigma_2 \in \Sigma$ such that

$$\forall y \neq x. \eta_y(\sigma_1(y)) = \eta_y(\sigma_2(y)) \land \neg \text{ATOM}_\rho(\llbracket e \rrbracket^\rho(\{\sigma_1, \sigma_2\})).$$

Since domains are partitioning, the nonatomicity requirement $\neg \text{ATOM}_\rho(\cdot)$ amounts to say that all abstract evaluations of $e$ with respect to $\rho$, starting from different values for $x$, may not be abstracted to the same abstract value (this is the crucial issue in $Ndep$); that is, $\sigma_1$ and $\sigma_2$ may lead to different abstract values for $e$. The next result shows that Atom-dep is an approximation of Ndep, since Ndep implies Atom-dep, meaning that Atom-dep may add false dependencies.

Proposition 6.6. Consider the abstractions $\rho \in uco(\wp(\Sigma))$ and $\overline{\eta} \in uco(\wp(\Sigma))^n$, where $n$ is the number of variables; that is, $\overline{\eta}$ is a tuple of properties. For every $e$ and $x$, $x \rightsquigarrow e$ implies $x \rightsquigarrow e$.

Proof. Suppose $x \rightsquigarrow e$; that is, there exist two states $\sigma_1$ and $\sigma_2$ such that (1) $\forall y \neq x. \eta_y(\sigma_1) =_y \eta_y(\sigma_2)$ and (2) $\rho(\llbracket e \rrbracket(\sigma_1)) \neq \rho(\llbracket e \rrbracket(\sigma_2))$. Then we have to prove $x \rightsquigarrow e$; that is, there exist $\sigma_1, \sigma_2 \in \Sigma$ such that (3) $\forall y \neq x. \eta_y(\sigma_1) =_y \eta_y(\sigma_2)$ and
(4) \(\lnot \text{Atom}_\rho([e]^{\rho}([\sigma_1, \sigma_2]))\). Since conditions (1) and (3) are the same, we have to prove that (2) implies (4).

Consider \(\sigma_1, \sigma_2 \in \Sigma\) satisfying (1) and such that \(\rho([e](\sigma_1)) \neq \rho([e](\sigma_2))\); then,
\[
\rho([e](\sigma_1)), \rho([e](\sigma_2)) \in \rho([e]([\sigma_1, \sigma_2])) \subseteq \rho([e](\rho([\sigma_1, \sigma_2]))) = [e]^{\rho}([\sigma_1, \sigma_2]).
\]
At this point, since \([e]^{\rho}([\sigma_1, \sigma_2])\) is greater than two different values of \(\rho\), then, by definition, it cannot be an atom of \(\rho\). ⊓⊔

A simple example where Atom-dep introduces false dependencies is the expression \(e = x \cdot x + x \cdot x + 1\), where \(\rho = \rho_{\text{PARSIGN}}\) (\(\overline{\eta}\) is irrelevant in this case). It is clear that the result of the expression is always positive and odd, so that \(\rho([e](\sigma)) = [\text{poseven}]\) for every \(\sigma\). Therefore, according to Definition 6.2, there is no abstract dependency of \(e\) on \(x\) with respect to \(\rho\). On the other hand, it is possible to take two concrete states \(\sigma_1\) and \(\sigma_2\) where \(x\) has different values and, in addition, such values are also different as regards the sign and/or the parity. For example, suppose \(\sigma_1(x) = 1\) and \(\sigma_2(x) = -4\): the abstract state \(\sigma^e\) that is the result of abstracting the set \(S = \{\sigma_1, \sigma_2\}\) is such that \(\sigma^e(x) = [\text{top}]\) since 1 and -4 have different parities and different signs. Computing the abstract semantics \([e]^{\rho}(\sigma^e)\) boils down to evaluate \(e\) at the abstract level with \([\text{top}]\) as the value for \(x\). Since \([\text{top}] \ast [\text{top}] + [\text{top}] \ast [\text{top}] + [\text{posodd}] = [\text{top}] + [\text{top}] + [\text{posodd}] = [\text{top}]^6\), then there is dependence of \(e\) on \(x\) according to Definition 6.5 since \([\text{top}]\) is obviously not an atom.

Starting from this new approximated notion, the idea is to provide an algorithm that overapproximates the set of variables relevant for a given expression \(e\). For simplicity, the abstraction observed on the output is supposed to be the same as the one on each input variable (i.e., each component of the tuple \(\overline{\eta}\) is exactly the \(\rho\) observed on the output, so that \(\overline{\eta}\) will be denoted as \(\overline{\rho}\) or even left implicit in the following). The idea is to start from an empty set of nonrelevant variables \(X\) and incrementally increase this set, adding all those variables that surely are not relevant. Finally, the complement of such set is returned, which is an overapproximation of relevant variables.

In order to check the dependency relation, we aim at checking whether a change in the values of a variable makes a difference in the evaluation of the expression. Dependencies are computed according to Atom-dep, in order to approximate Ndep. In a brute-force approach, Atom-dep would be verified by checking for each \(\sigma^e\) associating atomic values to variables we have in which \(\text{Atom}_\rho([e]^{\rho}(\sigma^e))\) is always the same atom.

**Example 6.7.** Let \(e\) be an expression involving variables \(x, y,\) and \(z\), and \(\rho = \rho_{\text{PARSIGN}}\). In principle, in order to compute the set of \(\rho\)-dependencies on \(e\), we must compute \([e]^{\rho}\) on every possible atomic value\(^7\) of \(x, y,\) and \(z\); that is, \([e]^{\rho}\) must be computed \(5^3 = 125\) times.

\(y\) is not relevant to \(e\) if, for any abstract values \(v_x, v_z \in \text{Atoms}(\rho)\), there exists an atomic abstract value \(v \in \text{Atoms}(\rho)\) such that \(\forall v \in \text{Atoms}(\rho). \ U = [e]^{\rho}([x \leftarrow v_x, y \leftarrow v, z \leftarrow v_z])\). This amounts to say that changing the value of \(y\) does not affect \(e\), since we require the same atom to be computed for each possible abstract value of \(y\) and \(z\). If the result is not atomic, it means that there may\(^8\) exist two different abstract results for different values of \(y\) and \(z\).

\(^7\)\([\text{posodd}]\) is the abstraction of the constant 1, which is positive and odd.
\(^8\)Remember that atoms in \(\rho\) are \([\text{zero}], [\text{poseven}], [\text{posodd}], [\text{negeven}],\) and \([\text{negodd}]\); since this is a partition of concrete values, we describe all concrete inputs by computing \([e]^{\rho}\) on atoms.

\(^8\)Since Atom-dep is just an overapproximation of semantic dependencies, it can be the case that such different results actually do not exist; on the other hand, atomicity does prove that no such results exist.
However, it is possible to be smarter:

- **Excluding states:** Consider dependencies of \( e \) in Example 6.7, computed at program point \( n \). Suppose \([e]^{\rho}\) (used as a tool to infer invariant properties, as discussed in Section 2.4) is able to infer, at point \( n \), that the abstract state \( \sigma_n^e \) such that \( \sigma_n^e(y) = [\text{posodd}] \) approximately approximates the value of variables at \( n \). Then, we only need to consider states of the form \( \{x \leftarrow v_x, y \leftarrow [\text{posodd}], z \leftarrow v_z\} \) as inputs for \([e]^{\rho}\) (now considered as the abstract computation of expressions, according to Definition 6.5) at \( n \).

- **Computing on nonatomic states:** Let \( E = \{[\text{poseven}], [\text{zero}], [\text{negeven}]\} \) and \( O = \{[\text{posodd}], [\text{negodd}]\} \). In this case,

\[
\forall \nu' \in E, \nu'' \in O. \ [e]^{\rho}((x \leftarrow v_x, y \leftarrow \nu', z \leftarrow \nu'')) \leq U
\]

is implied by the more general result

\[
[e]^{\rho}((x \leftarrow v_x, y \leftarrow [\text{even}], z \leftarrow [\text{odd}])) \leq U
\]

since \( E \) and \( O \) are partitions, respectively, of \([\text{even}]\) and \([\text{odd}]\), and \([e]^{\rho}\) is monotone: \( \sigma_1^e \leq \sigma_2^e \) implies \( \Rightarrow [e]^{\rho}(\sigma_1^e) \leq [e]^{\rho}(\sigma_2^e) \). This means that results obtained on \( \sigma^e \) can be used on \( \sigma_1^e \leq \sigma^\rho \).

In the following, we compute dependencies with respect to \( \sigma_1^\rho \), which is the abstract state computed at \( n \) as an invariant at that program point. In the worst case, no information about the different variables is available, and \( \sigma_1^\rho \) associates \([\text{top}]\) to all variables. Now, we aim at proving that \( e \) is independent from a set of variables \( X \); in order to prove this, we need to prove that evaluating \( e \) always yields the same atom in \( \rho \), independently from the value of variables in \( X \), possibly without trying all possible values. Our idea is to prove this atomicity, if it holds, by iteratively refining (Section 2.4) the starting abstract state \( \sigma_1^\rho \). The following example gives the intuition.

**Example 6.8.** Let \( e \equiv x * x + 1 \) and \( \rho = \rho_{\text{SIGN}} \). Let also \([e]^{\rho}\) follow the usual rules on * and +: \([\text{pos}] * [\text{pos}] = [\text{pos}], [\text{neg}] * [\text{neg}] = [\text{pos}], [\text{top}] * [\text{top}] = [\text{top}], [\text{pos}] + [\text{pos}] = [\text{pos}], [\text{top}] + [\text{pos}] = [\text{top}]\), and so forth. Suppose starting from a memory such that \( \{x \leftarrow [\text{top}]\} \); we observe that \([e]^{\rho}((x \leftarrow [\text{top}])) \in \text{ATOMS}(\rho_{\text{SIGN}}) \) cannot be proved by using these rules, since \([\text{top}] * [\text{top}] + [\text{pos}] = [\text{top}]\), which is not atomic, so that there may be a dependency. On the other hand, consider the possible refinements w.r.t. \( x \) in \( \rho_{\text{SIGN}} \), namely, \( \{x \leftarrow [\text{pos}], x \leftarrow [\text{zero}], x \leftarrow [\text{neg}]\} \}. These abstract values are a covering (Section 2.4) of \( \{x \leftarrow [\text{top}]\} \) since all possible values of \( x \) are accounted for by the refinements “\( x \) is positive,” “\( x \) is zero,” “\( x \) is negative.” Moreover, it turns out that

\[
[e]^{\rho}((x \leftarrow [\text{pos}])) = [e]^{\rho}((x \leftarrow [\text{zero}])) = [e]^{\rho}((x \leftarrow [\text{neg}])) = [\text{pos}].
\]

Thus, the positivity of \( x \) can be proved by using a kind of case-based reasoning that requires computing \( e \) three times (plus the initial one on \( \{x \leftarrow [\text{top}]\} \)) instead of just once.

Given an initial abstract state \( \sigma^\rho \), we define the set of all abstract states that “refine” abstract values on variables outside \( X \) to atomic values, while leaving \( X \) untouched.

\[
[\sigma^\rho|X] = \{ \sigma^\rho_y \mid \forall y \notin X. \sigma^\rho_y(y) \in \text{ATOMS}(\rho) \land \forall x \in X. \sigma^\rho_y(x) = \sigma^\rho(x) \}
\]

To prove \( e \) independent from \( x \), we need to prove \( \text{ATOM}_{\rho}([e]^{\rho}(\sigma^\rho_x)) \) for any \( \sigma^\rho_x \in [\sigma^\rho_n]|X \). This amounts to say that any variation in \( x \) (up to the invariant \( \sigma^\rho_n(x) \)) does not lead to an observable variation in \( e \), whenever all the other variables are fully specified.
as atoms. The following definition is based on the notion of covering introduced in Section 2.4 and identifies how an abstract state \( \sigma^\rho \) can be incrementally “refined,” that is, how it can be replaced by a set of states that is a covering, and where values of variables outside some set \( X \) take values that are the same as either the corresponding values in \( \sigma^\rho \) or their direct subvalues.

\[
X\text{-covering}(\sigma^\rho) = \left\{ \sigma_1^\rho, \ldots, \sigma_k^\rho \right\} \bigg| \left\{ \sigma_1^\rho, \ldots, \sigma_k^\rho \right\} \text{ is a covering and } \forall i. \forall x \in X. \sigma_i^\rho(x) = \sigma^\rho(x) \land \forall y \not\in X. \sigma_i^\rho(y) \leq \sigma^\rho(x) \right\},
\]

where \( \sigma_1^\rho(x) \leq \sigma_2^\rho(x) \) if and only if \( \sigma_1^\rho(x) = \sigma_2^\rho(x) \) or \( \sigma_1^\rho(x) \) is a direct subvalue of \( \sigma_2^\rho(x) \) in \( \rho \) (e.g., in \( \rho_{\text{ParSIGN}} \), \textit{posseven} is a direct subvalue of \textit{pos} but \textit{bot} is not because there is another abstract value between them). It is easy to see that, by repeatedly computing \( X\)-coverings starting from \( \sigma^\rho \), one ends up generating all abstract states in \( [\sigma^\rho]|_X \). This is what the algorithm presented in the following does incrementally.

**Example 6.9.** Consider \( \rho_{\text{ParSIGN}} \) and the abstract state \( \sigma^\rho = \{ x \leftarrow [\textit{pos}], y \leftarrow [\textit{odd}], z \leftarrow [\textit{top}] \} \): the following set of states is a \( \{z\}\)-covering of \( \sigma^\rho \).

\[
\left\{ \begin{array}{l}
\{ x \leftarrow [\textit{posseven}], y \leftarrow [\textit{odd}], z \leftarrow [\textit{top}] \}, \\
\{ x \leftarrow [\textit{posodd}], y \leftarrow [\textit{odd}], z \leftarrow [\textit{top}] \}
\end{array} \right.
\]

Another \( \{z\}\)-covering is the following, which refines both \( x \) and \( y \):

\[
\left\{ \begin{array}{l}
\{ x \leftarrow [\textit{posseven}], y \leftarrow [\textit{posodd}], z \leftarrow [\textit{top}] \}, \\
\{ x \leftarrow [\textit{posseven}], y \leftarrow [\textit{negodd}], z \leftarrow [\textit{top}] \}
\end{array} \right.
\]

\[
\left\{ \begin{array}{l}
\{ x \leftarrow [\textit{posodd}], y \leftarrow [\textit{posodd}], z \leftarrow [\textit{top}] \}, \\
\{ x \leftarrow [\textit{posodd}], y \leftarrow [\textit{negodd}], z \leftarrow [\textit{top}] \}
\end{array} \right.
\]

Moreover, given an expression \( e \) and a set of variables \( X \), we define

\[
\begin{align*}
\mathcal{A}_e^U(\sigma^\rho) & \quad \text{if and only if} \quad [e]^\rho(\sigma^\rho) = U \land \text{ATOM}_\rho(U) \lor \\
\exists(\sigma_1^\rho, \ldots, \sigma_k^\rho) \text{ covering of } \sigma^\rho \text{ such that } \forall i. \mathcal{A}_e^U(\sigma_i^\rho) \\
\mathcal{A}_e^*(\sigma^\rho, X) & \quad \text{if and only if} \quad \begin{cases} 
(1) \mathcal{A}_e^U(\sigma_i^\rho) \text{ for some atom } U \text{ or } \\
(2) \exists(\sigma_i^\rho)_{i \in \{1, k\}} \in X\text{-covering}(\sigma^\rho), \forall i. \mathcal{A}_e^*(\sigma_i^\rho, X).
\end{cases}
\end{align*}
\]

The predicate \( \mathcal{A}_e^U(\sigma^\rho) \) holds if the expression \( e \) can be proved to have an atomic value \( U \), either by direct computation of the abstract semantics or by case-based reasoning like in Example 6.8, relying on the notion of covering. On the other hand, \( \mathcal{A}_e^*(\sigma^\rho, X) \) holds if it can be proven that there is no dependency of \( e \) on variables in \( X \). This can be proven (1) directly, if \( e \) has an atomic value in \( \sigma^\rho \) (which implies that there is no dependency on variables in \( X \) or any other variables, or (2) recursively, by replacing \( \sigma^\rho \) with an \( X\)-covering and trying to prove the result for every refinement. Condition (2) relies on case-based reasoning (which is similar to the one outlined in Example 6.8, but not to be confused with it) and allows one to indirectly prove nondependency.

**Example 6.10.** Consider the expression \( e = 2^x y \), the \( \rho_{\text{ParSIGN}} \) domain, and the initial abstract state \( \sigma_0^\rho = \{ x \leftarrow [\textit{pos}], y \leftarrow [\textit{top}] \} \) (i.e., it is already known that \( x \) is positive). The predicate \( \mathcal{A}_e^*(\sigma_0^\rho, \{x\}) \) does not hold directly by condition (1) since the value of \( e \) in \( \sigma_0^\rho \) is obviously not atomic, and there does not exist any covering that allows one to prove \( \mathcal{A}_e^U(\sigma_0^\rho) \) for some \( U \). Then, \( \{x\}\)-coverings have to be searched for, which allows one to prove condition (2). One such \( \{x\}\)-covering is (the other \( \{x\}\)-coverings are not treated explicitly):

\[
\left\{ \begin{array}{l}
x \leftarrow [\textit{pos}], y \leftarrow [\textit{even}], \\
x \leftarrow [\textit{pos}], y \leftarrow [\textit{odd}]
\end{array} \right.
\]

Condition (2) would be proved if both \( \mathcal{A}_e^*(\{x \leftarrow [\textit{pos}], y \leftarrow [\textit{even}], \{x\}) \) and \( \mathcal{A}_e^*(\{x \leftarrow [\textit{pos}], y \leftarrow [\textit{odd}], \{x\})) \) hold. However, none of them hold: for example, the abstract
evaluation of \(\text{poseven} \times \text{pos} + \text{[even]}\) gives the nonatomic value \([\text{even}]\). Therefore, it is necessary to perform another recursive step: an \(\{x\}\)-covering of \(x \leftarrow \text{[pos]}, \ y \leftarrow \text{[even]}\) would be, for example,

\[
\{x \leftarrow \text{[pos]}, \ y \leftarrow \text{[poseven]}\}, \ \{x \leftarrow \text{[pos]}, \ y \leftarrow \text{[negeven]}\}, \ \{x \leftarrow \text{[pos]}, \ y \leftarrow \text{[zero]}\}.
\]

Unfortunately, it is still not possible to prove \(\tilde{\mathbb{A}}(\rho, X)\) for all three abstract states in the \(\{x\}\)-covering, and they cannot be further refined because values for \(y\) are already atomic. Therefore, nondependency of \(e\) on \(x\) cannot be proved. For the sake of simplicity, this example does not consider explicitly all \(\{x\}\)-coverings, but the result we give is true: \(\tilde{\mathbb{A}}(\rho, x)\) does not hold.

Example 6.11. Consider the same expression \(e = \text{2} \times x + y\), the \(\rho_{\text{ParSign}}\) domain, and the initial abstract state \(\sigma_n^* = \{x \leftarrow \text{[pos]}, y \leftarrow \text{[even]}\}\). In this case, it is possible to prove \(\tilde{\mathbb{A}}(\rho, x)\) because there is an \(\{x\}\)-covering:

\[
\{x \leftarrow \text{[pos]}, \ y \leftarrow \text{[poseven]}\}, \ \{x \leftarrow \text{[pos]}, \ y \leftarrow \text{[posodd]}\}.
\]

which satisfies condition (2). Therefore, if it is known that both \(x\) and \(y\) are positive, then it is possible to prove the independency of \(e\) from \(x\).

Proposition 6.12 (Soundness). Let \(\rho \in \omega(x \lor \forall)\), and let \(\vec{\rho}\) denote the tuple of \(\rho\), on each variable. If \(\tilde{\mathbb{A}}(\rho, X)\) can be proved, then there is no \(x \in X\) such that \(x \nvdash x \wedge \rho \wedge \phi \rightarrow e\).

Proof. By induction. Consider the definition of \(\tilde{\mathbb{A}}(\rho, x)\):

—If \(\tilde{\mathbb{A}}_e(\sigma_n^i)\) holds for some atomic abstract value \(\sigma\), then \(e\) depends on \(\sigma\) variables.

—Otherwise, let \(\sigma_1^i, \ldots, \sigma_k^i\) be an \(X\)-covering of \(\sigma_n^i\); by inductive hypothesis, for every \(i\), \(e\) can be proved not to depend on \(X\) if values are limited by \(\sigma_i^x\). Since these states are also a (regular) covering of \(\sigma_n^x\), and variables in \(X\) are never restricted in any of them (i.e., all values allowed by the initial \(\sigma_n^x\) are considered for every \(i\)), this means that \(e\) does not depend on \(X\).

The \textsc{findNdeps} algorithm depicted in Figure 14 starts by trying to prove \(\tilde{\mathbb{A}}(\rho, \text{vars}(e))\), where \text{vars}(e) are all the variables syntactically occurring in the expression. If it succeeds, then \(e\) depends on no variables. Otherwise, the variable set \(X\) is decreased nondeterministically (one element at a time, randomly) until some judgment \(\tilde{\mathbb{A}}(\rho, X)\) is proved.

The soundness of the algorithm follows easily from Proposition 6.12. On the other hand, the algorithm is also complete with respect to Atom-dep (which obviously does not mean that it exactly computes semantic abstract dependencies, due to the fact that Atom-dep is already an overapproximation of Ndep). In fact, exploring incrementally the space of \(X\)-coverings of \(\sigma_n^x\) is a way to consider all states in \([\sigma_n^x | X]\) without the need to generate all of them explicitly. Indeed, if nondependency from \(X\) can already
be proved for an abstract state belonging to an X-covering, then there is no need to consider all its refinements belonging to \([σ^o|X]\).

Importantly, the assertions \(A_e^U(σ^o)\) and \(A_e^U(σ^o,X)\) guarantee \(X \cup Y \not\sim \sim \sim \not\sim e\) not to hold, even if \(A_e^U(σ^o,X)\) cannot be directly proved. The final result of FINDDEPS is \(\{σ^1,σ^2\}\) is an X-covering of \(σ^o\) if it holds, even if \(A_e^U(σ^o,X)\) cannot be directly proved. The FINDDEPS algorithm may deal, in principle, with infinite abstract domains, in particular those with infinite descending chains, since nondependency results can be possibly proved without exploring the entire state space; in fact, as pointed out earlier, if \(A_e^U(σ^o)\) can be proved, then it is not needed to descend into the (possibly infinite) set of substates of \(σ^o\). This is not possible in the brute-force approach. It is also straightforward to add computational bounds in order to stop refining states if some amount of computational effort has been reached.

6.2.2. Dependency Erasure. The problem of computing abstract dependencies can be observed from another point of view: given \(e\) and a set \(X\) of variables, we may be interested in the most concrete \(ρ\) such that \(X \sim \sim \sim \not\sim e\) does not hold. This would be the most concrete observation guaranteeing the noninterference of the variables in \(X\) on the evaluation of \(e\) [Giacobazzi and Mastroeni 2004a; Mastroeni 2013]. This task can be accomplished by repeatedly simplifying an initial domain \(ρ_0\) in order to eliminate abstract values that are responsible for dependencies. As discussed before, in order not to have dependencies on \(X\), we should have \(A_e^U(σ^o,\{X\})\); that is, \(A_e^U(σ^o,\{X\})\) should hold for any \(σ^o \in [σ^o|X]\). If this does not hold for some abstract state \(σ^o\), then \(ρ\) is modified to make \(\llbracket e\rrbracket^{\prime}(σ^o)\) atomic.

We design a simple algorithm EDEP\((e, ρ_0, X)\) (Figure 15), which takes as input an expression, an initial uco, and a set of variables. The algorithm uses a queue in order to store abstract states; initially, the queue only contains \(σ^o_0\). When a state \(σ^o\) is extracted from the queue, EDep checks (line 6) if there exists \(u\) atomic such that \(A_e^U(σ^o)\). If it does not, then

— if it is not possible to further refine \(σ^o\) (guard at line 7), then the abstract domain must be simplified (see later), and the algorithm restarts from \(σ^o\) with the new domain;
—otherwise, $\sigma^\rho$ is refined according to the definition of $X$-covering (lines 13–15), and the process is repeated on all such refinements.

Each iteration of the loop corresponds to a value of the variable $\rho$. If the algorithm never reaches line 10 during a whole execution of a loop iteration (i.e., if, for every state $\sigma^\rho$, the condition $A^U(\sigma^\rho)$ can be proved at line 6 before $\sigma^\rho$ becomes “nonrefinable” at line 7), then $\rho$ is not simplified, which means that $A^U(\sigma^\rho, X)$ could be proved, thus demonstrating the nondependency of $e$ on $X$ with respect to $\rho$. In fact, the goal of lines 6 to 15 is exactly to decide if $A^U(\sigma^\rho, X)$, either directly (line 6) or considering $X$-coverings.

The simplifying operator ATOMIZE($\rho, \nu$) is a domain transformer and works by removing abstract values in order to obtain ATOM$_\nu$($\nu$). Formally,

$$\rho' = \text{ATOMIZE}(\rho, \nu) \overset{\text{def}}{=} \{ u \in \rho \mid \nu \cap u = \bot \lor \nu \leq u \}.$$

**Example 6.13.** Consider $\rho_{\text{ParSigh}}$; in order to make $[\text{even}]$ be atomic, it is necessary to remove abstract values that are not disjoint, greater than or equal to $[\text{even}]$. In this case, all values but $[\text{odd}]$ have to be removed from the abstract domain, so that $\text{ATOMIZE}(\rho_{\text{ParSigh}}, [\text{even}]) = \rho_{\text{Par}}$.

The final $\rho$ is an approximation of the most precise $\rho'$ s.t. $X \not\sim e \rightarrow e$ is false:

**Theorem 6.14.** $\rho \overset{\text{def}}{=} E_{\text{Deep}}(e, \rho_0, X)$ makes $e$ not Atom-dep on $X$, that is, for every $\sigma^\rho$ that is atomic on $\text{vars}(e) \setminus X$, $[e]^\rho(\sigma^\rho)$ is atomic. Moreover, by Proposition 6.6, $e$ is not dependent on $X$ with respect to Definition 6.2.

**Proof.** The algorithm halts if, in processing $\sigma^\rho_n$, $\rho$ is not changed. Processing $\sigma^\rho_n$ involves computing $[e]^\rho$ on substates when required, in order to prove the atomicity property on every concrete state represented by $\sigma^\rho_n$ (to this end, we exploit monotonicity of $[e]^\rho$ on states). This is precisely obtained if every state is removed from the queue before any modification to $\rho$ occurs. □

**Example 6.15.** Consider the expression $e = 2 \times x + y$, the $\rho_{\text{ParSigh}}$ domain, and the set of variables $\{x\}$. It is clear that $e$ does depend on $x$ with respect to $\rho_{\text{ParSigh}}$. Suppose that the algorithm starts from $\sigma^\rho_0 = \{ x \leftarrow [\text{top}], y \leftarrow [\text{top}] \}$. Since $[e]^\rho(\sigma^\rho_0) = [\text{top}]$ and the condition at line 7 is false, an $\{x\}$-covering of $\sigma^\rho_0$ is to be considered (line 15):

$$\{ x \leftarrow [\text{top}], y \leftarrow [\text{even}] \}.$$ For each abstract state in the $\{x\}$-covering, the same test is performed, but the value of the expression is still nonatomic. The computation goes on without computing atomic values for $e$, until the guard at line 7 becomes true (e.g., when considering $\{ x \leftarrow [\text{top}], y \leftarrow [\text{poseven}] \}$). Since the abstract value computed in this state is $[\text{even}]$, the procedure ATOMIZE($\rho_{\text{ParSigh}}, [\text{even}]$) is called at line 10, and the new candidate abstract domain $\rho_{\text{Par}}$ is computed (Example 6.13). Next, the queue is reset to its initial value, and another iteration of the loop begins. Following the same reasoning, it is easy to see that the dependency no longer holds in the new domain $\rho_{\text{Par}}$, so that the second iteration of the loop does not simplify the domain anymore, and $\rho_{\text{Par}}$ is finally returned.

On the practical side, the loss of precision in abstract computations may lead one to remove more abstract values than strictly necessary from the semantic point of view. It is important to note that EDeep works as long as $A_u$ can be computed on the initial domain (in this case, no problems arise in subsequent computations, since the “complexity” of $\rho$ can only decrease). This can possibly happen even if $\rho_0$ is infinite (see the end of Section 6.2.1). Moreover, unlike FINDNEEPS, there is no reasonable trivial counterpart, so any brute-force approach would be definitely impractical.
7. THE QUEST FOR ABSTRACT SLICES

This section introduces an algorithm for computing conditioned abstract slices, based on abstract dependencies and the notion of agreement between states. As explained in Section 7.4, this way to compute slices relies on a priori knowledge of the properties that will be of interest for the analysis. In most cases, the majority of the abstract domains to be taken into account are quite simple (e.g., nullity). On the other hand, Section 6.2 presents a general way to compute abstract dependencies on more complex domains. Indeed, those algorithms can be used here, and their complexity is acceptable if small domains are dealt with.

The slices obtained by following this methodology will have the standard form of backward slicing; predicates or conditions $\beta$ on states are supported, in the style of conditioned slicing [Canfora et al. 1998]. In the rules, the judgment $\sigma \models \beta$ means that the state $\sigma$ satisfies $\beta$. It must be pointed out that a predicate $\beta$ at a certain program point may include user-provided information but also knowledge about program variables that is obtained by performing static analysis on the program (see the discussion in Section 2.4 and the use of $\sigma_{\rho}^n$ in Section 6.2). For example, after a $x:=\text{new } C()$ statement, judgments like “$x$ is not null,” “$x$ is not cyclic,” “$x$ is not sharing with $y$,” or “$x$ is not reaching $y$” could be provided, depending on the kind of static analyses available (nullity, sharing, cyclicity, etc.).

A way to decide which statements have to be included in an abstract slice consists of two main steps:

—For each statement $s$, a specific static analysis algorithm provides information about the relevant data after that statement (below, the agreement), according to the slicing criterion and the program code.
—If the execution of $s$ does not affect its corresponding agreement (i.e., some condition on states that must hold after $s$), then $s$ can be removed from the slice.

Example 7.1. Consider the following code fragment, and suppose that the slicing criterion is the nullity of $x$ at the end:

```plaintext
21 . . .
22 y . f := exp;
23 x := y;
```

The field update on $y$ (line 22) can be removed from the slice because

— the question about the nullity of $x$ after line 23 is equivalent to the question about the nullity of $y$ after line 22; and
— the field update at line 22 does not affect the nullity of $y$.

The rest of this section formalizes how these two main steps are carried out.

7.1. The Logic for Propagating Agreements

This section describes how agreements are defined and propagated via a system of logical rules: the $G$-system. Hoare-style triples [Hoare 1969] are used for this purpose, in the spirit of the weakest precondition calculus [Dijkstra 1975].

Definition 7.2. An agreement $G$ is a set of conditions $[\bar{x}:\rho]$ where each uco $\rho$ involves a sequence of variables $\bar{x}$ (most usually, just one variable), and all conditions involve

---

9A version of this system was introduced in previous work [Zanardini 2008] as the $A$-system. However, there are many differences between both systems, mainly due to the changes in the language under study (e.g., variables are taken into account here instead of a more involved notion of pointer expression).
mutually disjoint sets of variables. Two states $\sigma_1$ and $\sigma_2$ are said to agree on $\mathcal{G}$, written $\mathcal{G}(\sigma_1, \sigma_2)$, if and only if, for every $[x::\rho]$ in $\mathcal{G}$, $\rho(\sigma_1(x)) = \rho(\sigma_2(x))$, where notation is abused by taking $\sigma_i(x)$ as the sequence of values of variables $x$ in $\sigma_i$, and $\rho(\sigma_i(x))$ as the application of $\rho$ (which could be a relational domain) to the elements of such a sequence.

Agreements are easily found to form a lattice, and a partial order $\sqsubseteq$ can be defined: $\mathcal{G} \leq \mathcal{G}''$ if and only if, for every $\sigma_1, \sigma_2$ such that $\mathcal{G}'(\sigma_1, \sigma_2)$, then $\mathcal{G}''(\sigma_1, \sigma_2)$. Moreover, an intersection operator is induced by the partial order: $\mathcal{G}' \cap \mathcal{G}''$ is the greatest agreement that is less than or equal to both.

In the following, $\mathcal{G}(x)$ will be the uco corresponding to the condition $[x::\rho]$ in $\mathcal{G}$, or $\rho_T$ if no condition on $x$ belongs to $\mathcal{G}$. For the sake of simplicity, the discussion will be limited to domains each involving one single variable. In this case, ordering amounts to the following: $\mathcal{G}_1 \subseteq \mathcal{G}_2$ if $\forall x. \mathcal{G}_1(x) \subseteq \mathcal{G}_2(x)$, where $\mathcal{G}_1(x)$ is the comparison on the precision of abstract domains, meaning that $\mathcal{G}_1(x)$ is more precise than $\mathcal{G}_2(x)$.

**Example 7.3.** Let $\sigma_1 = \{n \leftarrow 2, i \leftarrow 3, x \leftarrow \text{null}\}$ and $\sigma_2 = \{n \leftarrow 0, i \leftarrow 4, x \leftarrow \text{null}\}$ be two states. Then, they agree on $\mathcal{G} = \{[(n)::\rho_{\text{par}}], [(x)::\rho_{\text{null}}]\}$ since

- $n$ has the same parity in both states;
- $x$ is $\text{null}$ (therefore, it has “the same nullity”) in both states; and
- there is no condition on $i$.

On the other hand, these states do not agree on $\mathcal{G}' = \{[(i)::\rho_{\text{par}}], [(x)::\rho_{\text{null}}]\}$ because $i$ is odd in $\sigma_1$ whereas it is even in $\sigma_2$.

In a triple $\{\mathcal{G}\}$ s $\{\mathcal{G}'\}$, the precondition $\mathcal{G}$ is the weakest agreement on two states before a statement $s$ such that the agreement specified by the postcondition $\mathcal{G}'$ holds after the statement. Predicates on states can be used, so that triples are, actually, 4-tuples that only take into account a subset of the states. Formally, the 4-tuple (or augmented triple) $\{\mathcal{G}\}^{\beta} s \{\mathcal{G}'\}$ (where the true predicate is often omitted) holds if, for every $\sigma_1$ and $\sigma_2$,

$$\sigma_1 \models \beta \quad \land \quad \sigma_2 \models \beta \quad \land \quad \mathcal{G}(\sigma_1, \sigma_2) \Rightarrow \mathcal{G}'([s](\sigma_1), [s](\sigma_2)).$$

The rules of the $\mathcal{G}$-system are shown in Figure 16. The transformed predicate $s(\beta)$ is one that is guaranteed to hold after a statement $s$, given that $\beta$ holds before, in the style of strongest postcondition calculus [Dijkstra and Scholten 1990; Bijlsma and Nederpelt 1998]. For example, if $\beta = x \geq 0$, then the condition $s(\beta) = x \geq 1 \land y \neq \text{null}$ certainly holds after $s = x := x + 1; y := \text{new} C(0)$. The way predicates are transformed is outside the scope of this article; however, the cited works introduced calculi for computing such strongest postconditions. In the absence of such tools, true is always a consistent choice (precision, but not soundness, may be affected since the set of states to be considered grows larger, and to prove useful results could become harder).

The $\mathcal{G}$-system is tightly related to narrow noninterference [Giacobazzi and Mastroeni 2004a, 2004b]. These works define a similar system of rules, the N-rules, for assertions $[\eta]s(\eta')$, where $\eta$ and $\eta'$ are basically the (tuples of) abstract domains corresponding to, respectively, $\mathcal{G}$ and $\mathcal{G}'$. The systems differ in that:

- the use of pointers requires the rules for assignment to account for sharing, while $\mathcal{G}$-rules only work on integers;
- in the present approach, domains are supposed to be partitioning, so that there is no need to include explicitly the $\Pi$ operator (Section 2.3);
- the $\mathcal{G}$-system does not distinguish between public and private since this notion is not relevant in slicing;
Fig. 16. The G-system.

—the rule for conditional is not included in the N-system; indeed, this is quite a tricky rule, and, in general, expressing a conditional with loops and using the rule N6 for loops results in inferring less precise assertions; and
—in the N-system, predicates \( \beta \) on program states are not supported.

In the following, each rule of the G-system is discussed. Importantly, this rule system relies on the computation of property preservation. We rely on a rule system that soundly computes whether executing a statement affects properties of some variables: the judgment

\[
\text{PP}^\beta(\mathcal{G}, s)
\]

can only be obtained by using those rules if it is possible to prove that executing \( s \) in a state \( \sigma \) satisfying the condition \( \beta \) results in a final state \( \sigma' \) that is equal to the initial \( \sigma \) with respect to the agreement \( \mathcal{G} \), that is, \( \mathcal{G}(\sigma, \sigma') \). In other words, the
statement is equivalent to skip with respect to the properties of interest. Note that, here, the agreement is not used to compare two states at the same program point; rather, it takes as input the states before and after executing a statement. The rule system for proving property preservation is explained after introducing the g-system (Section 7.1.7).

7.1.1. Rule G-pp. This rule makes the relation between property preservation and the g-system more clear. The triple \( \{G\}_s \{G\} \) amounts to say that two executions agree after \( s \), provided they agree before on the same \( G \). On the other hand, the preservation of \( G \) on \( s \) means that any state before \( s \) agrees on \( G \) with the corresponding state after \( s \). Property preservation is a stronger requirement than the mere propagation \( \{G\}_s \{G\} \) of agreements, so that this rule is sound. In fact, if \( G(\sigma_1, \sigma_2) \) and both \( G(\sigma_1, [s](\sigma_1)) \) and \( G(\sigma_2, [s](\sigma_2)) \) hold, then \( G([s](\text{state}_1), [s](\sigma_2)) \) follows, which is equivalent, by definition, to \( \{G\}_s \{G\} \).

Example 7.4. Let the parity of \( x \) be the property of interest. In this case, \( x := x+1 \) does not preserve the parity of \( x \), but two initial states agreeing on \( [x::\rho_{\text{par}}] \) lead to final states that still agree on it. Therefore, \( [x::\rho_{\text{par}}] x := x+1 \) \( [x::\rho_{\text{par}}] \) holds. On the other hand, \( x := x+2 \) also satisfies a stronger requirement: that \( \rho_{\text{par}}(x) \) does not change. Therefore, besides having \( [x::\rho_{\text{par}}] x := x+2 \) \( [x::\rho_{\text{par}}] \), the judgment \( \text{pp}(x::\rho_{\text{par}}, x := x+2) \) is also true.

7.1.2. Rules G-skip, G-concat, G-id, G-sub. The G-skip rule describes no-op. The assertion holds for every \( G \) and \( \beta \) since \( \llbracket \text{skip} \rrbracket(\sigma) = \sigma \).

G-concat is also easy: soundness holds by transitivity (note also the use of \( s(\beta) \) to propagate conditions of states). Rule G-id can be used when nothing else can be proved; it always holds because execution is deterministic, so that two executions starting from two states that are equal\(^{10}\) on all variables will end in a pair of states agreeing on any abstraction. Note that, as pointed out in Section 2.1, read statements are supposed only to appear at the beginning of a program; therefore, a compound statement never contains any read, and the propagation of agreements can be done safely until reaching the first (actually, the last occurrence in the code) read statement (in the presence of read statements, there is no need to propagate agreements until the very first line of the code).

Finally, in G-sub, remember that \( \sqsubseteq \) is the partial order on agreements.

7.1.3. Rule G-assign. This rule means that, given a statement \( x := e \), any agreement \( G \) that satisfies the two conditions of the previous part of the rule is a sound precondition for the postcondition \( G' \). The conditions are that (1) given two states that agree on \( G \), the computed results for the expression \( e \) in both states are abstracted by \( G'(x) \) to the same abstract value, and (2) \( G \) is as precise as \( G' \) on all variables but \( x \). The first condition is represented in terms of Definition 6.5, and the superscript \( \beta \) indicates that only states satisfying \( \beta \) have to be considered. Such a condition can be easily shown to imply the formula
\[
F = \forall \sigma_1 \models \beta, \sigma_2 \models \beta . \; (G(\sigma_1, \sigma_2) \Rightarrow (G'(x)(\llbracket e \rrbracket(\sigma_1)) = (G'(x)(\llbracket e \rrbracket(\sigma_2)))) .
\]
Note that, by Proposition 6.6, the absence of abstract dependencies w.r.t. Definition 6.5 (Atom-dep) implies the absence of abstract dependencies w.r.t. Definition 6.2 (Ndep), which in turn implies \( F \). This clarifies the relation between abstract dependencies and the computation of a slice. Obviously, the second condition guarantees that, for all variables that are not updated by the assignment, the agreement required by \( G \) still holds when \( G' \) is considered.

\(^{10}\)The notion of equality on references and objects is recalled in Example 2.2.
7.1.4. Rule G-FASSIGN. This rule accounts for the modification of the data structure pointed to by a variable by means of a field update. In the following, that a variable is affected means that the data structure pointed to by it is updated. Given a field update on a variable \( x \), some other variables (i.e., the data structures pointed to by them) could be affected. There exists a well-known static analysis that tries to detect which variables point to a data structure that is updated by a field update on \( x \): this analysis is known as sharing analysis [Secci and Spoto 2005; Zanardini 2015] and usually comes as possible-sharing analysis, where the set of variables that could be affected by a field update is overapproximated. Moreover, aliasing analysis [Hind 2001] can be used in order to compute the set of variables pointing exactly (and directly) to the same location as \( x \); in this case, definite-aliasing analysis makes sense, which underapproximates the set of variables that certainly alias with \( x \). According to the result of these analyses, reference variables can be partitioned in three categories: (1) variables that certainly alias with \( x \), so that they can be guaranteed to be updated in their field \( f \); (2) variables not in (1) that could share with \( x \), so that they could be affected by the update in many ways; and (3) variables that certainly do not share with \( x \), so that they are unaffected by the update. Let \( \text{SH}(x) \) be the set of variables possibly sharing with \( x \) before the update, and \( \text{DAL}(x) \) be the set of variables definitely aliasing with \( x \). In the absence of a definite-aliasing analysis, \( \text{DAL}(x) \) can be safely taken as \( \{x\} \).

\[ \text{Example 7.5.} \] Consider the following code fragment:

\begin{verbatim}
10 if (...) then { y.f := x; } else { y.f := z; }
11 w := x;
12 x.g := e;
\end{verbatim}

Suppose that, initially, no variable is sharing with any other variable (i.e., there is no overlapping between data structures referred by different variables), and that the truth value of the Boolean guard cannot be determined statically, so that both branches of the conditional statement have to be considered as possible executions. In this case, the sharing and aliasing information before line 12 is as follows:

\[ \text{SH}(x) = \{x, y, w\} \quad \text{DAL}(x) = \{x, w\}. \]

Note that every variable is aliasing with itself (none of them is null), \( w \) is certainly aliasing with \( x \) because of line 11, and \( y \) is possibly sharing with \( x \) (the actual sharing depends on the value of the guard).

The G-FASSIGN rule comes with three preconditions. Precondition \((*)\) only applies to variables in category (1): those definitely aliasing with \( x \). Precondition \((***)\) applies to category (2), while precondition \((***)\) applies to categories (2) and (3).

\[-(*)\] requires that updating the field \( f \) of the location pointed to by \( x \) (and all variables definitely aliasing with it) leads to an agreement on \( G' \), provided that the initial states agree on \( G \).

\[-(**)\] is similar, but states that the agreements must hold for every sequence of field selectors \( \bar{g} \), possibly the empty sequence. This is needed since, given that some \( y \) may share with \( x \), it cannot be known which fields of \( y \) will be updated, and how. In practice, only the sequences of field selectors that are compatible with the

\[11\text{In the following, the notation \( \langle .f_1.f_2.\ldots.f_n \rangle \) (starting with a dot is intentional) will be used to represents sequences of field selectors.}\]
class hierarchy of the program under study have to be considered, as shown in Example 7.6.

—(***) applies to variables that could be unaffected by the update, that is, variables in categories (2) and (3) (note that the conditions $y \in \text{SH}(x)$ in (***) and $y \notin \text{DAL}(x)$ in (***) are not mutually exclusive, so that variables in category (2) satisfy both). The relation between $\mathcal{G}$ and $\mathcal{G}'$ is clear in this case, as agreement on $\mathcal{G}$ must entail agreement on $\mathcal{G}'$.

Example 7.6 (Infeasible Sequences of Field Selectors). In a Java-like language, a sequence like $\langle .g.f \rangle$ is not compatible with the following class hierarchy since the class of $g$ is $D$, whereas $f$ is declared in $C$:

```
class C { D f; D g; }
```

```java
class D { D h; }
```

Example 7.7. Consider the statement $s \equiv x.f := y$. Let $\beta = \text{true}$ and the agreement $\mathcal{G}'$ after $s$ be $\{[x::\rho_{\text{NULL}}], [z::\rho_{\text{NULL}}]\}$. Let also $\text{SH}(x)$ before $s$ be $\{x\}$. In this case, an agreement $\mathcal{G}$ that satisfies the judgment

$$\{\mathcal{G}\}_{\text{true}} x.f := y \{\mathcal{G}'\}$$

can be the same $\{[x::\rho_{\text{NULL}}], [z::\rho_{\text{NULL}}]\}$ because

—$z$ is unaffected by the update (it belongs to category (3)), so that $\mathcal{G}(z)$ must be at least as precise as $\mathcal{G}'(z)$;

—category (2) contains no variables; and

—category (1) only includes $x$ itself, and the nullity of $x$ is clearly unaffected by the update.

One may think that the universal quantification on field sequences in precondition (***) results in an unacceptable loss of precision. Indeed, to require that all possible updates to possibly sharing variables preserve the desired agreements seems to be too strict. However, there are a number of things to be considered:

—A closer look to the rule shows that there is no easier way to account for sharing if traditional sharing analysis is used.

—Example 7.7 shows that it is still possible to get meaningful results on domains working on pointer variables.

—The state of the art in static analysis of object-oriented languages indicates that abstract domains on pointers are likely to be quite simple ($\rho_{\text{NULL}}$ being one of them).

—There is recent work [Zanardini 2015] introducing a more precise, field-sensitive sharing analysis that computes how variables share: this analysis is able to detect which fields are or are not involved in paths in the heap converging from two variables to a shared location. In order to keep the discussion as simple as possible, the definition of $\mathcal{G}$-$\text{FASSIGN}$ given in Figure 16 uses traditional sharing analysis. However, the impact of field-sensitive sharing analysis is discussed in Section 7.4, where a refined version of $\mathcal{G}$-$\text{FASSIGN}$ is given.

The domain of cyclicity introduced in Section 2.2 represents information about data structures in the heap, not only program variables. In this sense, field updates have to be regarded as potentially affecting the propagation of agreements.

Example 7.8. Let $s$ be, again, the statement $x.f := y$, and $\mathcal{G}'$ be $\{[x::\rho_{\text{CYC}}]\}$; that is, the interest is on the cyclicity of the data structure pointed to by $x$ after executing $s$. Suppose also that $x$ and $y$ are certainly not sharing before $s$, and that an object whose type is compatible with $x$ has two reference fields $f$ and $g$. Then, the agreement
Fig. 17. How two executions agreeing on $\{x::\rho_{\text{CYC}}, y::\rho_{\text{CYC}}\}$ do not agree on $\{x::\rho_{\text{CYC}}\}$ after $x.f:=y$.

$G = \{x::\rho_{\text{CYC}}, y::\rho_{\text{CYC}}\}$ is not a correct precondition for the Hoare tuple to hold, since there can be two states $\sigma_1$ and $\sigma_2$ (Figure 17) such that:

— the data structure corresponding to $y$ is acyclic, and equal in both states (therefore, there is an agreement on the cyclicity of $y$);
— the data structure corresponding to $x$ is cyclic in both $\sigma_1$ and $\sigma_2$, but (1) in $\sigma_1$ there is only a cycle originating from the location bound to $x.f$, and (2) in $\sigma_2$ there is a cycle originating from $x.f$ and another one originating from $x.g$.

In this case, $G(\sigma_1, \sigma_2)$ holds but the resulting final states $\sigma'_1$ and $\sigma'_2$ do not agree on $G'$ since $x$ is acyclic in $\sigma'_1$ (the only cycle has been broken) while it is still cyclic in $\sigma'_2$. This behavior is captured by $G$-\textsc{fassign} because condition $(\ast)$ applied to $x$ itself (which, by hypothesis, is the only variable in $\text{DAL}(x)$) does not hold, so that the augmented triple cannot be proven. On the other hand, $\{x::\rho_{\text{ID}}, y::\rho_{\text{CYC}}\}$ would be a correct precondition.

7.1.5. Rules $G$-IF1 and $G$-IF2. In a conditional if $(b) \ s_t \ else \ s_f$ there are two possibilities. Rule $G$-IF1 states that an input agreement that induces the output one whichever path is taken is a sound precondition. Here, the assertion $\{G\}^\beta \ s_t \ o s_f \ (G')$ means that

$$\forall \sigma_1, \sigma_2. \ G(\sigma_1, \sigma_2) \land \sigma_1 \models \beta \land \sigma_2 \models \beta \Rightarrow G'(\llbracket s_t \rrbracket(\sigma_1), \llbracket s_t \rrbracket(\sigma_2), \llbracket s_f \rrbracket(\sigma_1), \llbracket s_f \rrbracket(\sigma_2)),$$

where the judgment $G'(\ldots, \ldots, \ldots)$ means that all four states agree on $G'$. This rule requires $G'$ to hold on the output state independently from the value of $b$. Soundness is easy (note that the previous assertion implies $\{G\}^\beta \ s_t \ o \ s_f \ (G')$).

Note that such a $G$ can always be found (in the worst case, it assigns the identity upper closure operator $\rho_{\text{ID}}$ to each variable, so that two states agree only if they are exactly equal). However, sometimes it can be more convenient to exploit information about $b$. In such cases, $G$-IF2 can be applied, which means that the initial agreement $G_t \cap G_f$ is strong enough to verify the final one, provided the same branch is taken in both executions, as $G_b$ requires. In fact, $G_b$ is built from $b$ and separates states according to its value:

$$G_b(\sigma_1, \sigma_2) \ Leftrightarrow \llbracket b \rrbracket(\sigma_1) = \llbracket b \rrbracket(\sigma_2).$$
The rule means that, whenever two states agree on the branch to be executed, and the triples on both branches hold, the whole triple holds as well.

**Lemma 7.9 (Soundness of g-if2).** If $\sigma_1$ and $\sigma_2$ both satisfy $\beta$ and agree on $G_b \cap G_t \cap G_f$, then the corresponding output states $\sigma_1'$ and $\sigma_2'$ agree on $G'$ under the hypotheses of the rule.

**Proof.** By hypothesis, the same branch is taken in both cases. Conditions $\beta \wedge b$ and $\beta \wedge \neg b$ are consistent since $s_t$ ($s_f$, respectively) can only be executed when $b$ is true (false, respectively). The agreement on $G$ holds in both paths, so that the entire assertion is correct. \qed

Note that the rule to be chosen for the conditional depends on the precision of the outcome: g-if2 can be a good choice if (1) it can be applied and (2) the result is “better” than the one obtained by g-if1. The second condition amounts to say that, given the same final agreement $G'$, the initial agreement obtained by using g-if2 is weaker (i.e., it is more likely that two states agree on it) than the one obtained by using g-if1.

**Example 7.10.** Consider the code fragment

\[
\text{if } (x > 0) \{ x := x + 1; \} \text{ else } \{ x := x - 1; \}
\]

and let $G = \{ [x::P] \}$ be the agreement after the statement; that is, the relevant property is the sign of $x$. The rule g-if2 is able to compute the same $G$ as the input agreement because

— the triple $\{ G \}^{\beta \wedge b} s_t \{ G \}$ holds since the condition $\beta \wedge b$ guarantees that $x$ is positive, and two states that agree on the sign before the increment will still agree after it (if $x$ is positive in both states, then it will remain positive in both);

— similarly, the triple $\{ G \}^{\beta \wedge \neg b} s_f \{ G \}$ also holds (if $x$ is 0 in both states, then it will be negative in both; and if it is negative in both, it will remain negative in both);

— $G_b$ is less precise than $G$ (i.e., $G \subseteq G_b$ since the latter only separates numbers into positive and nonpositive), so that the input agreement $G_b \cap G \cap G$ is equal to $G$.

On the other hand, g-if1 is not able to compute the same input agreement because the precondition $\{ G \}^\beta s_t \circ s_f \{ G \}$ of the rule does not hold. In fact, consider two states $\sigma_1 = \{ x \leftarrow 1 \}$ and $\sigma_2 = \{ x \leftarrow 2 \}$: they agree on the sign of $x$, but $([s_t](\sigma_2))(x)$ and $([s_f](\sigma_1))(x)$ have different signs (the first is zero while the second is positive).

**7.1.6. Rule g-while.** The meaning of the rule for loops can be understood by discussing its soundness: if $\beta$ is preserved after any iteration of the body, and the agreement that is preserved by the body guarantees the same number of iterations in both executions (i.e., it is more precise than $G_b$), then such an agreement is preserved through the entire loop.

**Lemma 7.11 (Soundness of g-while).** Let $\sigma_1^0$ and $\sigma_2^0$ satisfy $\beta$ and agree on $G \cap G_b$. Then, given $\sigma_1' = [\text{while } (b) \ s_n] \sigma_1^0$, the result $(G \cap G_b)(\sigma_1', \sigma_2')$ holds.

**Proof.** Let $\sigma_i^{n+1} = [s_n](\sigma_i^n)$ for every $n \geq 0$; that is, $\sigma_i^{n+1}$ refers to the situation after the $(n+1)$-th iteration. There are two cases, depending on the value of the Boolean condition:

— $[b](\sigma_i^n) = [b](\sigma_2^n) = \text{false}$: in this case, the body is not executed, and the result holds trivially.

— $[b](\sigma_1^n) = [b](\sigma_2^n) = \text{true}$: in this case, $\sigma^n \models \beta \wedge b$. By the hypothesis of the rule, $\sigma_1^{n+1}$ and $\sigma_2^{n+1}$ agree on $G \cap G_b$, and $\beta$ still holds since $s_n(\beta) \Rightarrow \beta$. 

ACM Transactions on Computational Logic, Vol. 18, No. 1, Article 7, Publication date: February 2017.
Note that \( G_b \) guarantees the same number of iterations in both executions; consequently, for a terminating loop (nontermination is not considered), \( \sigma_{1}^{k} \) and \( \sigma_{2}^{k} \) will fall in the first case (false guard) after the same number \( k \) of iterations. These states are exactly \( \sigma_{1}^{'} \) and \( \sigma_{2}^{'} \) and agree on \( G \sqcap G_b \) after the loop.

**Theorem 7.12 (G-Soundness).** Let \( s \) be a statement, \( G' \) be required after \( s \), \( \beta \) be a predicate, and \( p \) be the program point before \( s \). Let also \( G \) be an agreement computed before \( s \) by means of the \( G \)-system. Let \( \tau_{1} \) and \( \tau_{2} \) be two trajectories, and the states \( \sigma_{1} \in \tau_{1}[p] \) and \( \sigma_{2} \in \tau_{2}[p] \) satisfy \( G(\sigma_{1}, \sigma_{2}) \) and \( \beta \). Then, the condition \( G'(\sigma_{1}^{'}, \sigma_{2}^{'}) \) holds, where \( \sigma_{1}^{'} = [s](\sigma_{1}) \).

**Proof.** Easy from Lemmas 7.9 and 7.11, and the discussion explaining each rule (especially, G-FASSIGN). \( \Box \)

### 7.1.7. The PP-System

Property preservation can be proved by means of a rule system, the PP-system (Figure 18). Most rules are straightforward or very similar to G-system rules and characterize when executing a certain statement preserves the properties represented by an agreement \( G \) (i.e., for every variable \( x \), the property/ucco \( G(x) \) is preserved). For example, rule PP-ASSIGN allows one to prove that a certain agreement is preserved when the initial value of \( x \) cannot be distinguished from the value of the expression (i.e., the new value of \( x \)), when it comes to the property \( G(x) \). In rule PP-FASSIGN, a mechanism similar to G-FASSIGN is used: definite aliasing and possible sharing can be used to identify which variables are affected by the field update. As for G-FASSIGN, an optimization based on field-sensitive sharing analysis (Section 7.4.2) can be introduced, which makes it easier to prove property preservation on field updates.

Indeed, a number of optimizations can be applied to the PP-system (e.g., one could think that property preservation on \( s_{1} ; s_{2} \) does not require the preservation of the same properties on both statements separately). However, this rule system is not the central part of this article, and Figure 18 is just a sensible way to infer property preservation.
7.2. Agreements and Slicing Criteria

This approach to compute abstract slices follows the standard conditioned, non-iteration-count form of backward slicing. Therefore, a slicing criterion \( C \) takes the form \((\mathcal{I}, \mathcal{X}, \{n\} \times \mathbb{N}, \text{false}, \mathcal{A})\), where \( n \) is the last program point, and \( \mathcal{A} \) is a sequence of ucos assigning a property to each variable in \( \mathcal{X} \). The initial agreement can be easily computed from the criterion and is such that \( \mathcal{G} = \{ (x::A_x) \mid x \in \mathcal{X} \} \), where \( A_x \) is the element of \( \mathcal{A} \) corresponding to \( x \). This shows that there is a close relation between this specific kind of slicing criteria and agreements, and in the following, these concepts will be used somehow interchangeably in informal parts. It will be shown that this makes sense; that is, criteria and agreements define tightly related notions. The next definition defines the correctness of an abstract slice of this kind, where the slicing criterion is intentionally confused with an agreement.

**Definition 7.13 (Abstract Slicing Condition).** Let \( P^* \) be the slice of \( P \) with respect to an agreement (criterion) \( \mathcal{G} \). In order for \( P^* \) to be correct, \( [P](\sigma) \) and \( [P^*](\sigma) \) must agree on \( \mathcal{G} \) for every initial \( \sigma : \mathcal{G}([P](\sigma)), [P^*](\sigma)) \).

7.3. Erasing Statements

The main purpose of the \( \mathcal{g} \)-system is to propagate a final agreement backward through the program code, in order to have a specific agreement attached to each statement.\(^{12}\) This is done as follows: a program can be seen as a sequence \( \mathcal{S} = s_0; \ldots; s_k \) of \( k + 1 \) statements, where each \( s_i \) can be either a simple statement (skip, assignment, field update) or a compound one (conditional or loop), containing one (the loop body) or two (the branches of the conditional) sequences of statements (either simple or compound, recursively). The way to derive an agreement for every statement in the program is depicted in the pseudocode of Figure 19. The procedure \( \text{labelSequence} \) takes as input:

1. A sequence of statements (in the first call, it is the whole program code\(^{13}\));
2. A pair of agreements: (2.a) the first one, \( \mathcal{G}_{in} \), refers to the beginning of the sequence and, in the first call, is such that the abstraction on each variable is \( \rho_{in} \), and (2.b) the second one, \( \mathcal{G}_{out} \), is the desired final agreement, which corresponds to the slicing criterion as discussed in Section 7.2; and
3. A predicate on states that is supposed to hold at the beginning of the sequence.

\( \text{labelSequence} \) goes backward through the program code inferring, for each \( s_i \), an agreement \( \mathcal{G}_i \) that corresponds to the program point after \( s_i \), \( \mathcal{G}_k \) will be the same \( \mathcal{G}_{out} \), whereas, for each \( i \), \( \mathcal{G}_{i-1} \) will be inferred by using the \( \mathcal{g} \)-system: more specifically, it is (ideally, the best) a precondition such that the tuple \( \{ G_{i-1} \}_{i-1}^{\beta_i} s_i \{ G_i \} \) holds. Note that, since the initial \( \mathcal{G}_0 \) is the identity on all variables, \( \{ G_0 \}_{i-1}^{\beta_i} s_i \{ G_i \} \) trivially holds for every \( i \) (execution is deterministic); however, the \( \mathcal{G}_{in} \) argument plays an important role when dealing with loop statements.

Statements inside compound statements (e.g., assignments contained in the branch of a conditional) are also labeled with agreements: this is done by calling \( \text{labelSequence} \) recursively. Note that, in this case, if \( \vec{s}_t \) and \( \vec{s}_f \) are, respectively, the sequences corresponding to the “then” and “else” branch of a conditional statement \( s_i \), and then \( \text{labelSequence} \) is called with second argument \( (\mathcal{G}_{in}, \mathcal{G}_i) \); this is so because the state does

\(^{12}\) Here, a statement is not only a piece of code but also a position (program point) in the program, so that no two statements are equal, even if they are syntactically identical. For the sake of readability, the program point is left implicit.

\(^{13}\) Strictly speaking, the sequence of statements is the program without the initial sequence of \text{read} statements: remember that \text{read} statements are basically meant to provide the input.
procedure labelSequence ('s_1;...;s_k ', (G_in, G_out), β) {
  G_0 = G_in;
  G_k = G_out;
  every β_i is \( \overline{s}(β) \) where \( \overline{s} = 's_1;...;s_i ' \);
  // (remember the transformed predicate \text{s}(β))
  for i = k downto 1 {
    if (s_i is a conditional) {
      let b be the guard;
      let \( \overline{s}_l \) and \( \overline{s}_f \) be its branches;
      call labelSequence (\( \overline{s}_l \), (G_0, G_i), β_{i-1} \land b);
      call labelSequence (\( \overline{s}_f \), (G_0, G_i), β_{i-1} \land \neg b);
      G_{i-1} is such that \{G_{i-1}\}_{ß_{i-1}} s_i \{G_i\};
    }
    else if (s_i is a loop) {
      G_{i-1} is an agreement such that
      - \( G_{i-1} \subseteq G_i \cap G_b \) // (see rule G-WHILE)
      - \{G_{i-1}\}_{ß_{i-1}} s_i \{G_i\}_{ß_i}  
      let b be the guard;
      let \( \overline{s}_l \) be the loop body;
      call labelSequence (\( \overline{s}_l \), (G_{i-1}, G_{i-1}), β_{i-1} \land b);
    }
    else { // non–compound statement
      G_{i-1} is such that \{G_{i-1}\}_{ß_{i-1}} s_i \{G_i\};
    }
  }
}

Fig. 19. Labeling program code with agreements by using the g-system.

not change when control goes from the end of a branch to the statement immediately after \( s_i \).

The treatment of loops follows closely the definition of the rule G-WHILE. In the augmented triple at line 16, \( G_{i-1} \) appears before and after the statement; this is consistent with the rule. The condition \( G_{i-1} \subseteq G_i \cap G_b \) guarantees that the \{\( G_{i-1} \}_{ß_{i-1}} s_i \{G_i\}_{ß_i} \) can be proven by applying the g-SUB rule. Moreover, the recursive call on the body \( \overline{s}_l \) at line 19 has \( (G_{i-1}, G_{i-1}) \) as its second argument. In this recursive call, the value of the formal parameter \( G_{in} \) will no longer be, in general, the identity on all variables. The code fragment dealing with loops (lines 14–19) relies on an agreement \( G_{i-1} \) satisfying a number of conditions (line 15, line 16, and the success of the recursive call at line 19). It is important to note that such an agreement always exists, so that the execution of the pseudoalgorithm never stops: the \( G_{i-1} \) satisfying all conditions is the identity on all variables.

The following proposition specifies when it is correct to slice out a statement from a program.

**Proposition 7.14.** Consider a program \( P \). Suppose that the judgment \( \text{pv}^\beta(G_s, s) \) can be proved, where \( G_s \) and \( β_s \) are computed by labelSequence, and \( β \) refers to the program point before \( s \). In this case, let \( P' \) be the program \( P \) where \( s \) has been replaced by \text{skip}, and \( σ \) be an initial state. Then, the following holds: upon termination (nontermination is not considered), \( [P](σ) \) agrees with \( [P'](σ) \) with respect to the agreement \( G_{out} \) corresponding to the desired slicing criterion.

**Proof.** Let \( τ \) and \( τ' \) be two trajectories coming from executing, respectively, \( P \) and \( P' \) from the initial state \( σ \). Let \( σ_{in}[j] \) and \( σ'_{in}[j] \) be the states of \( τ \) and \( τ' \), respectively, when control reaches the program point before \( s \) (or \text{skip}, in the case of \( P' \)) for the \( j \)th time, and \( σ_{out}[j] \) and \( σ'_{out}[j] \) be their corresponding states after \( s \) (or \text{skip}). If \( s \) is not contained in any loop, then \( j \) can only be 1, and the proof is trivial. Otherwise, it can be any number up to some \( k_s \) (a nonnegative number).
Fig. 20. Graphical representation of some aspects of Proposition 7.14. \( \tau \) and \( \tau' \) are the executions of \( P \) and \( P' \), respectively. Labeled horizontal double arrows indicate the (provable) agreements between both traces at different program points: (a) at the beginning of the loop body; (b) before \( s \) or \( \text{skip} \); (c) after \( s \) or \( \text{skip} \); and (d) at the end of the body. Moreover, both \( s \) and \( \text{skip} \) preserve \( G_s \).

If \( k_s = 0 \), then the loop is never executed on the input \( \sigma \), and the proof follows trivially. Otherwise, \( \sigma_{\text{in}[1]} \) and \( \sigma_{\text{in}'}[1] \) are identical because both executions have been going exactly through the same statements; as a consequence, they certainly agree on \( G_s \). On the other hand, \( \sigma_{\text{out}[1]} \) and \( \sigma_{\text{out}'}[1] \) are in general not identical, but they still agree on \( G_s \) by the following:

1. \( G_s(\sigma_{\text{in}[1]}, \sigma_{\text{in}'}[1]) \) (they are identical)
2. \( G_s(\sigma_{\text{in}[1]}, \sigma_{\text{out}[1]}) \) (semantics of \( \text{skip} \))
3. \( G_s(\sigma_{\text{in}[1]}, \sigma_{\text{out}[1]}) \) (property preservation on \( s \))
4. \( G_s(\sigma_{\text{out}[1]}, \sigma_{\text{out}'}[1]) \) (transitivity applied to (1), (2), and (3))

Due to how the program is labeled with agreements by labelSequence, the agreement \( G_s(\sigma_{\text{out}[1]}, \sigma_{\text{out}'}[1]) \) after the first iteration implies an agreement of both executions at the end of the loop body with respect to the agreement \( G_{\text{end}} \) labeling that program point. This also means that both executions will still agree on \( G_s \) at the beginning of the next (second) iteration, and (again, by construction) they will also agree on \( G_s \) when \( s \) is reached for the second time. This mechanism can be repeated until \( k_s \) is reached, and it is easy to realize that the agreement on \( G_s \) at the last iteration implies the final agreement on \( G_{\text{out}} \) (see Figure 20 for a picture).

The crux of this reasoning is that, by construction, agreements labeling each program point imply that, when a statement can be removed from the body, \( P \) and \( P' \) will execute the loop body the same number of times. In general, this does not mean that every loop will be executed the same number of times (e.g., a property-preserving loop could be even sliced out completely as in Example 7.17).

Example 7.15. Consider the following program:

```plaintext
23 while (y>0) { // \{[x)::\rho_{\text{PAR}}, [x)::\rho_{\text{SIGN}}]\}
24   if (...) { x := x+1; }
25   ... 
26   x := x - 2; // \{[x)::\rho_{\text{PAR}}, [x)::\rho_{\text{SIGN}}]\}
27 } // \{[x)::\rho_{\text{PAR}}]\}
```

Given the final agreement, the loop cannot be entirely sliced out because parity may change at line 24. The agreement after lines 23 and 26 is the \( G_{i-1} \) introduced at line 14.
of labelSequence; it clearly satisfies the condition at line 15 of the same algorithm, and we suppose that it also satisfies the augmented triple of line 16. Then, removing the statement at line 26 is irrelevant to the slicing criterion since it does not modify the parity of \( x \) nor the sign of \( y \).

This proposition can be used in order to remove all statements for which such a property preservation judgment can be proved. In general, the slice is computed by replacing all statements \( s \) such that \( \text{pp}^\beta(G, s) \) holds by \text{skip} or, equivalently, removing all of them from the original code. It is easy to observe that this corresponds exactly to the notion of backward abstract slicing given in Section 4.

**Example 7.16.** Consider the following code, and let the final nullity of \( x \) be the property of interest (corresponding to the agreement \( \{ \{ x::\rho_{\text{null}} \} \} \) after line 15):

```plaintext
8 ... 
9 n := n*2; \quad // \{ \{ \tau::\rho_{\text{zero}} \} \}
10 C x := new C(); \quad // \{ \{ \tau::\rho_{\text{zero}} \} \}
11 if (n=0) { \quad // \{ \{ \text{final result on this branch always null} \} \}
12 \quad x := null; \quad // \{ \{ x::\rho_{\text{null}} \} \}
13 } else { \quad // \{ \{ \text{final result on this branch never null} \} \}
14 \quad x := new C(); \quad // \{ \{ x::\rho_{\text{null}} \} \}
15 }
```

Each agreement on the right-hand side is the label after the statement at the same line. Both lines 9 and 10 can be removed from the slice because:

— the final nullity of \( x \) only depends on the equality of \( n \) to 0 before line 11, which is captured by the \( \rho_{\text{zero}} \) domain;
— line 10 does not affect \( n \); and
— multiplying a number by 2 preserves the property of being equal to 0.

Note that, although agreements after lines 11 and 13 are both empty (no matter how the state is, \( x \) will be null after line 12 and nonnull after line 14), the one after line 10 is not because of rule g-if2.

**Example 7.17.** Consider again the code of Example 4.1. Starting from the final agreement \( \{ \{ \text{list}::\rho_{\text{cycle}} \} \} \), the code is annotated as follows:

```plaintext
33 y := null; \quad // \{ \{ \text{list}::\rho_{\text{cycle}} \} \}
34 x := list; \quad // \{ \{ \text{list}::\rho_{\text{cycle}}, x::\rho_{\text{cycle}} \} \}
35 while (pos>0) { \quad // \{ \{ \text{list}::\rho_{\text{cycle}}, x::\rho_{\text{cycle}} \} \}
36 \quad y := x; \quad // \{ \{ \text{list}::\rho_{\text{cycle}}, x::\rho_{\text{cycle}} \} \}
37 \quad x := x.next; \quad // \{ \{ \text{list}::\rho_{\text{cycle}}, x::\rho_{\text{cycle}} \} \}
38 \quad pos := pos - 1;
39 }
40 z := new Node(elem); \quad // \{ \{ \text{list}::\rho_{\text{cycle}}, x::\rho_{\text{cycle}} \} \}
41 z.next := x; \quad // \{ \{ \text{list}::\rho_{\text{cycle}}, z::\rho_{\text{cycle}}, x::\rho_{\text{cycle}} \} \}
42 if (y = null) { \quad // \{ \{ \text{list}::\rho_{\text{cycle}}, z::\rho_{\text{cycle}}, x::\rho_{\text{cycle}} \} \}
43 \quad list := z; \quad // \{ \{ \text{list}::\rho_{\text{cycle}} \} \}
44 } else { \quad // \{ \{ \text{list}::\rho_{\text{cycle}}, z::\rho_{\text{cycle}}, x::\rho_{\text{cycle}} \} \}
45 \quad y.next := z; \quad // \{ \{ \text{list}::\rho_{\text{cycle}} \} \}
46 }
```

To prove the necessary tuples, it is important to note that next is the only reference field selector in the class Node, so that any cycle has to traverse it. Moreover, to prove that the cyclicity of list after line 46 is related to the cyclicity of list, z, and x before
that line needs nontrivial reasoning about data structures; concretely, it is necessary to have some reachability analysis [Genaim and Zanardi 2013]\(^\text{14}\) capable to detect that there is no path from \(z\) to \(y\) (otherwise, a cycle could be created by \(y.\text{next} = z\)). This information should be available as \(\beta\)\(^\text{15}\) and allows one to say that \(y\) and \(\text{list}\) (note that \(\text{list}\) is affected because it is sharing with \(y\)) are cyclic after line 46 if and only if \(z\) or \(x\) or \(\text{list}\) were before that line (actually, a closer inspection of the code reveals that \(z\), \(x\), and \(\text{list}\) are either all cyclic or all acyclic before line 46). Variable \(y\) does not appear in the agreement at line 45 because next is the only field, so that any cycle possibly reachable from it before line 46 is no longer reachable after that same line.

The situation is similar at line 42: the cyclicity of \(\text{list}\) and \(x\) is not affected by the field update since \(z\) is provably not sharing with any of them before line 42. On the other hand, the cyclicity of \(z\) depends on \(x\) and \(\text{list}\) because any cycle reachable from \(z\) after line 42 must be also reachable from the other two variables.

Moreover, the agreement does not change in the loop body at lines 36 to 40 because (1) data structures are not modified, (2) \(\text{list}\) is not affected in any way, and (3) the value of \(x\) changes, but its cyclicity does not (by executing \(x := x.\text{next}\), there is no way to make unreachable a cycle that was reachable before, or the other way around). Importantly, this also means that the cyclicity of \(\text{list}\) and \(x\) is preserved by the loop, so that it can be safely removed from the slice (the preserved property is the same as the agreement after line 40).

On the other hand, the conditional statement at lines 43 to 47 cannot be removed directly because it is not possible to prove that the cyclicity of \(\text{list}\) is preserved through it (actually, it is not). In order to prove that the whole code between lines 34 and 47 preserves the cyclicity of \(\text{list}\), a kind of case-based reasoning could be used: (1) the initially acyclicity of \(\text{list}\) implies its final acyclicity, and (2) the initially cyclicity of \(\text{list}\) implies its final cyclicity. Both these results can be proved by standard static analysis techniques [Genaim and Zanardi 2013].

Needless to say, slices could be suboptimal (e.g., the requirement about executing loops the same number of times need not be satisfied by any correct slice). It is not difficult to see that, if a “concrete” slicing criterion would be considered instead of an abstract one, then agreements would only be allowed to contain conditions \([x::\rho_m]\) for a certain set of variables. The labelSequence would work exactly the same way, with an important difference: when trying to compute the precondition of a tuple \(\{\_\}^\beta G\), the possible outcome could, again, contain only conditions \([x::\rho_m]\).

**Example 7.18.** In Example 7.16, suppose the final agreement is \([x::\rho_m]\), corresponding to a concrete slicing criterion interested in (the exact value of) \(x\). In this case, the agreement after line 10 could only be \([n::\rho_m]\), since (1) {} would not be correct, and (2) no other abstract domain can appear in agreements. This way, line 9 could not be sliced out since it does not preserve the \(\rho_m\) property of \(n\).

Semantically, abstract slices are in general smaller than concrete slices. This is also the case of Example 7.18. Clearly, this does not imply that every abstract slicing algorithm would remove more statements than every concrete slicing algorithm.

---

\(^{14}\)Note that pair-sharing-based cyclicity analysis [Rossignoli and Spoto 2006] is not enough since \(y\) and \(z\) are sharing before line 46.

\(^{15}\)See Section 2.4 and the beginning of this section for a discussion about \(\beta\) and similar examples of statically inferred information about states.
7.4. Practical Issues and Optimizations

This section discusses how the analysis can realistically deal with the computation of abstract dependencies and the propagation of agreements, and an optimization based on recent work on sharing.

7.4.1. Agreements and ucors. One of the major challenges of the whole approach is how agreements are propagated backward through the code as precisely as possible, that is, being able to detect that agreement on some ucors before \( s \) implies agreement on (possibly) other ucors after \( s \).

Ideally, given \( s, \beta, \) and \( \mathcal{G}' \), the \( \mathcal{G} \)-system should find the best \( \mathcal{G} \) such that \( \{\mathcal{G}\}^\beta \models \{\mathcal{G}'\} \). However, it is clearly unrealistic to imagine that the static analyzer will always be able to find the best ucors without going into severe scalability (even decidability) issues: in general, there exist infinite possible choices for an input agreement satisfying the augmented triple. In practice, an implementation of this slicing algorithm will be equipped with a library of ucors among which the satisfaction of augmented triples can be checked. The wider the library is, the more precise the results. Rule of the \( \mathcal{G} \)-system and the \( \mathcal{G} \)-system can be specialized with respect to the ucors at hand.

Example 7.19. Consider this code already presented in Example 7.16, where the slicing criterion is the final nullity of \( x \):

```plaintext
8  ...  
9   n := n * 2; 
10  C x := new C(); 
11  if (n=0) { 
12       x := null; 
13   } else { 
14       x := new C(); 
15   }
```

In order to be able to remove lines 9 and 10 from the slice, an analyzer has to “know” that, after merging the result of both branches, the uco \( \rho_{\text{ZERO}} \) precisely describes the agreement before line 11. In other words, the analyzer must know both \( \rho_{\text{ZERO}} \) and \( \rho_{\text{NULL}} \) in order to be able to manipulate information about them. Moreover, given a library of ucors, rules can be optimized for some recurrent programming patterns like guards (\( n=0 \)) or (\( x=null \)), or statements \( m:=0 \).

It is clear that designing an analyzer that is able to deal with all possible ucors is infeasible. However, the combination of some simple numeric or reference domains like the one described in this article would already lead to meaningful results.

7.4.2. Use of Field-Sensitive Sharing. As already mentioned, field-sensitive sharing analysis [Zanardini 2015] is able to keep track of fields that are involved in converging paths from two variables to a common location in the heap (the shared location). A propositional formula is attached to each pair of variables and each program point, and specifies the fields involved in every pair of converging paths reaching a common location. For example, if the formula \( \neg \chi_x \land \mathcal{F} \) is attached to a pair of variables \( (x, y) \) at a certain program point \( n \) (written \( S_n(x, y) = \neg \chi_x \land \mathcal{F} \)), this means that the analysis was able to detect that, for every two paths \( \pi_1 \) and \( \pi_2 \) in the heap starting from \( x \) and \( y \), respectively, and both ending in the same (shared) location,

\(-\pi_1\) certainly does not traverse field \( f \), as dictated by \( \neg \chi_x \) (arrows from top-left to bottom-right refer to paths from \( x \), i.e., the first variable in the pair under study); and
—π₂ certainly traverses g, as prescribed by $\mathcal{K}$ (arrows from top-right to bottom-left refers to paths from y).

In the presence of such an analysis, two kinds of improvements can be potentially obtained when analyzing a field update:

—The number of variables that can be actually affected by an update is, in general, reduced since it is possible to guarantee that some (traditionally) sharing variables will not be affected.
—Even for variables that are (still) possibly sharing with x, the set of field sequences $\bar{g}$ to be considered can be substantially smaller.

Example 7.20. Suppose that field-sensitive sharing analysis is able to guarantee the following at a program point before the field update $x.f := e$:

—The formula $\neg \chi \land \mathcal{K}$ correctly describes the sharing between x and y.
—The formula $\mathcal{K}$ correctly describes the sharing between x and z.

According to the traditional notion of sharing, both y and z may share with x. However:

—the assignment $x.f := e$ provably does not affect y because no path from x traversing f will reach a location that is also reachable from y; and
—when considering all the possible field sequences starting from z, only those containing h have to be considered.

The rule g-fassign can be refined by using field-sensitive sharing, as follows. Let $n$ be the program point before the field update.

—In precondition ($\ast\ast$), the only sharing variables that have to be dealt with are those for which it cannot be proved that they are unaffected by the update; this can be done by defining a new set $SH_f(x)$ of variables that are possibly sharing with x in such a way that some path from x to a shared location could traverse f:

$SH_f(x) = \{ y \mid S_n(x, y) \not\models \neg \chi \}.$

This means that y is considered as potentially affected by the update when the propositional formula describing how it shares with x does not entail that paths from x to shared locations do not traverse f.

—In the same precondition ($\ast\ast$), the universal quantification on field sequences can be restricted to those compatible with field-sensitive information. More formally, a field sequence $\langle f_1, f_2, ..., f_n \rangle$ has to be considered only if it is possible that a path from y traversing exactly those fields ends in a shared location or, equivalently, if the set $\{ \chi, \mathcal{K}, \mathcal{K}, ..., \mathcal{K} \}$ is a model of $S_n(x, y)$. That such a set is a model of $S_n(x, y)$ is equivalent to saying that the field-sensitive information is compatible with the existence of a pair of paths $\pi_1$ and $\pi_2$ such that (1) $\pi_1$ starts from x; (2) $\pi_2$ starts from y; (3) $\pi_1$ only traverses f; (4) $\pi_2$ traverses all and only the fields $f_1, f_2, ..., f_n$; and (5) both paths end in the same shared location. Let $\textit{fields}(\bar{g})$ be the set of fields contained in the field sequence $\bar{g}$. Then, the previous condition can be written as

$SH_{seq}^{x,y}(f, \bar{g}) = \bigwedge_{p \in X} p \land \bigwedge_{q \notin X} \neg q \models S_n(x, y),$

where $X = \{ \chi \} \cup \{ g \mid g \in \textit{fields}(\bar{g}) \}$, and $q \notin X$ means that q is any proposition $\chi$, or $\mathcal{K}$ (for some field h) not included in X.
The refined \text{G-Fassign} rule, called \text{G-Fassign2}, comes to be

\[
\forall y \in \text{DAL}(x) . \forall \sigma_1 \models \beta, \sigma_2 \models \beta. \\
\sigma_1 \models [y. f \leftarrow [e](\sigma_1)], \sigma_2 \models [y. f \leftarrow [e](\sigma_2)].
\]

\text{G-Fassign2.}

\[\beta \]

7.5. Comparison with Related Algorithms

The Tukra Abstract Program Slicing Tool [Halder and Cortesi 2012] implements the computation of a dependence condition graph [Halder and Cortesi 2013] for performing abstract slicing. To use a program dependence graph is somehow alternative to the computation of agreements. As far as the authors make it possible to understand, Tukra only deals with numerical values, and it is not clear which properties are supported (i.e., which is the “library” of ucos mentioned in Section 7.4.1).

Moreover, the authors of that tool point out that their approach is able to exclude some dependencies that were not ruled out in previous work introducing abstract dependencies [Mastroeni and Zanardini 2008]. However, they do not consider that a rule system for computing agreements (essentially, the \text{G-system} described in the present article) was introduced [Zanardini 2008] before Tukra was developed and does not suffer from the limitations they describe.

8. RELATED WORK

The formal framework referred to in this article [Binkley et al. 2006a] is not the only attempt to provide a unified mathematical framework from program slicing. In Ward and Zedan [2007], the authors have precisely this aim. In this work, the authors unify different approaches to program slicing by defining a particular semantic relation, based on the weakest precondition semantics, called \text{semirefinement} such that, given a program \( P \), the possible slices are all the programs that are semirefinements of \( P \).

In this framework, different forms of slicing are modeled as program transformations. Hence, a program \( Q \) is a slice of \( P \) if the transformation of \( Q \) (corresponding to the particular form of slicing to compute) is a semirefinement of the same transformation of \( P \). This approach is extremely interesting, but does not really allow one to compare the different forms of slicing, a feature that we consider fundamental for introducing the new abstract forms of slicing as generalizations of the existing ones. It may surely deserve further research to study whether abstract slicing could also be modeled in this framework.

As far as the relation between slicing and dependencies is concerned, there are at least two works that are related with our ideas in different ways. One of the first works aiming at formalizing a semantic approach to dependency, leading to a semantic computation of slicing, is the \text{information-flow logic} by Amtoft and Banerjee [2007]. This logic allows us to formally derive, by structural induction, the set of all the \text{independencies} among variables. In Figure 21, the original notation proposed by the authors is used, where \([x \times y]\) is to be read as “the current value of \( x \) is independent of the initial value of \( y \)” and holds if, for each pair of \text{initial} states that agree on all the variables but \( y \), the corresponding \text{current} states agree on \( x \). Hence, \( T^\# \) stands for sets of independencies, and \( G \) is a set of variables representing the \text{context}, that is, (a superset of) the variables on which at least one test surrounding the statements depends on.
In our aim of defining slicing in terms of dependencies, the first thing we have to observe in this logic is that it always computes (in)dependencies from the initial values of variables. This makes its use for slicing not so straightforward, since it loses the local dependency between statements. Consider, for example, the program fragment:

\[ P = w := x + 1; y := w + 2; z := y + 3 \]

At the end of this program, we know that \( z \) only depends on the initial value of \( x \), but, by using the logic in Figure 21, we lose the trace of (in)dependencies that, in this case, would involve all three assignments. As a matter of fact, this logic is more suitable for forward slicing, which is the one considered by Amtoft and Banerjee [2007], since it fixes the criterion on the input. In the trivial example given previously, if we consider as criterion the input of \( x \), then we obtain that all the statements depend on \( x \). Therefore, any slice of the original program contains all statements [Amtoft and Banerjee 2007]. In the logic, more explicitly, this notion of dependency is used for characterizing the set of independencies holding during the execution of a program.

Another, more recent, approach to slicing by means of dependencies is Danicic et al. [2011]. In this work, the authors propose new definitions of control dependencies: nontermination sensitive and insensitive. These new semantic notions of dependencies are then used for computing more precise standard slices. It could be surely interesting to study the semantic relation between their notion of dependencies and the ones we propose in this article.

Finally, a related algorithm for computing abstract slices has already been discussed in Section 7.5. It is necessary to point out that the agreement-based approach to abstract slicing [Zanardini 2008] was introduced before the Tukra tool.

9. CONCLUSION AND FUTURE WORK

The present article formally defines the notion of abstract program slicing, a general form of slicing where properties of data are observed instead of their exact value. A formal framework is introduced where the different forms of abstract slicing can be compared; moreover, traditional, nonabstract forms of slicing are also included in the framework, allowing one to prove that nonabstract slicing is a special case of abstract slicing where no abstraction on data is performed.

Algorithms for computing abstract dependencies and program slices are given. Future work includes an implementation of this analysis for an object-oriented programming language where properties may refer either to numerical or to reference values (to the best of our knowledge, existing tools only deal with integer variables).
On the other hand, we observed that the provided notion of abstract dependency is not suitable for slicing computation by using PDGs. We believe that it is possible to further generalize the notion of abstract dependencies allowing one to characterize a recursive algorithm able to track backward both the variables that affect the criterion and the abstract properties of these variables affecting the abstract criterion.

Another interesting line of research is to understand how other approaches to slicing can be extended in order to include abstract slicing. As noted before, it would be interesting to study whether it is possible to model abstract slicing as a program transformation, allowing us to define also abstract slicing in terms of semirefinement [Ward and Zedan 2007]. Another, more algorithmic, interesting approach is the one proposed in Barros et al. [2010], where weakest precondition and strongest postcondition semantics are combined in a new more precise algorithm for standard slicing. It could be very interesting to understand whether this approach could be extended in order to cope also with the computation of abstract forms of slicing.

APPENDIX
A. THE FORMAL FRAMEWORK OF PROGRAM SLICING

In this section, we first provide a better intuition of the differences between the forms of slicing introduced in Section 3.1 by means of examples and then we recall the main notions introduced in Binkley et al. [2006a, 2006b] that have been generalized in this article in the abstract form.

Different Forms of Slicing: Some Examples

Example A.1. Consider the program on the left in Figure 22. Suppose that the execution starts with an initial value 2 for \( n \) (written \( n \leftarrow 2 \)). States are denoted as \((m^k, \mu)\), where \( m \) is the program point of the executed statement, \( k \) is its current iteration (i.e., the statement at \( m \) is being executed for the \( k \)th time in the loop unrolling), and \( \mu \) is the actual memory, represented by a list of pairings \( x \leftarrow v \). In the picture, \( m^k \) is depicted in the first box, while the memory is depicted in the remaining boxes, one for each variable. A program state that is not executed in a trace is depicted by overwriting a gray cross on each box (program point and variables). The execution trajectory is the following:

Consider now the code in the center, whose execution trajectory is the following:

precisely the same path on the statements that are in both programs (the only difference is the execution, in the original program of the statement at point 6, erased in the “candidate” slice); hence, it is a slice according both to the standard and the KL form, namely, w.r.t. \( C = (\{n \leftarrow 2\}, \{x\}, \{(11, N)\}, \psi) \) for both possible values of \( \psi \).

Suppose now we are interested in an IC form of slice, and the value of \( x \) at the second iteration of the program point 3. In this case, the program on the right is not a dynamic
The value of \( x \) in the original program is 18, while it is undefined in the candidate slice. In other words, this program is not a slice of the program on the left w.r.t. the criterion \( C = ([n \leftarrow 2], \{x\}, \{(3, \{2\}\}), \psi) \) (\( \psi \) may be both true or false).

Finally, let us consider the execution of the program on the right in Figure 22:

This last program is a standard dynamic slice since the final value of the variable of interest \( x \) is the same, but it is not a slice in the \( KL \) form, since in this last program the statement at program point 10 is executed, while in the original program it is not executed. Namely, it is a slice w.r.t. the criterion \( C = ([n \leftarrow 2], \{x\}, \{(11, \mathbb{N})\}, \psi) \) only for \( \psi = false \).

The following example shows the difference between standard and \( KL \) forms of static slicing.

**Example A.2.** Consider the program \( P \) on the left of Figure 23. Suppose static slicing is considered; that is, all the possible initial memories are taken into account. Given an input \( v \), the state trajectory is:

Consider now the code on the right: its state trajectory is:

Fig. 22. Programs of Example A.1.

Fig. 23. Programs of Example A.2.
Consider the standard form of static slicing interested in \( x \) at program point \( 8^{1} \), that is, \( C = ((n \leftarrow \mathbb{N}), \{x\}, \{(8, \{1\})\}, \text{false}) \). Then, the program on the right is a slice of \( P \) w.r.t. \( C \), since in \( 8^{1} \) the value of \( x \) is \( v + 2 \) in both cases. On the other hand, if we consider the KL form, \( C = ((n \leftarrow \mathbb{N}), \{x\}, \{(8, \{1\})\}, \text{true}) \), then the program is not anymore a slice of \( P \) since there is a program point, \( 7^{1} \), that is not reached in \( P \).

The Unified Equivalence

The first step for defining the formal framework is to define an equivalence relation between programs, determining when a program is a slice of another. First, a restricted memory is obtained from a memory by restricting its domain to a set of variables. More formally, the restriction of \( \mu \) with respect to a set of variables \( \mathcal{X} \) is defined as \( \mu \upharpoonright \mathcal{X} \) such that \( (\mu \upharpoonright \mathcal{X})(x) \) is equal to \( \mu(x) \) if \( x \in \mathcal{X} \), and undefined otherwise. This restriction is used to project the trace semantics only on those points of interest where we have to check the correspondence between the original program and the candidate slice.

The trajectory projection operator modifies a state trajectory by removing all those states that do not contain occurrences of program points that are relevant for the slicing criterion.

**Definition A.3 (Trajectory Projection [Binkley et al. 2006a]).** Let \( C = (\mathcal{I}, \mathcal{X}, \mathcal{O}, \psi) \) and \( \mathcal{L} \subseteq \mathbb{L} \) such that \( \mathcal{L} \neq \emptyset \) if \( \psi = \text{true} \), and \( \mathcal{L} = \emptyset \) if \( \psi = \text{false} \). For any \( n \in \mathbb{L}, k \in \mathbb{N}, \mu \in \mathbb{M} \), we define the function \( \text{Proj}_{0} \) as:

\[
\text{Proj}_{0}(\mathcal{X}, \mathcal{O}, \mathcal{L})(n^{k}, \mu) \overset{\text{def}}{=} \begin{cases} 
(\langle n^{k}, \mu \upharpoonright \mathcal{X} \rangle) & \text{if } (\exists(n, K) \in \mathcal{O}. k \in K) \\
(\langle \varepsilon, \bot \rangle) & \text{if } (\exists(n, K) \in \mathcal{O}. k \in K) \text{ and } n \in \mathcal{L} \\
\varepsilon & \text{otherwise},
\end{cases}
\]

where \( \varepsilon \) is the empty sequence. The trace projection \( \text{Proj} \) is the extension of \( \text{Proj}_{0} \) to sequences (\( \circ \) is sequence concatenation):

\[
\text{Proj}(\mathcal{X}, \mathcal{O}, \mathcal{L})(\langle n^{k_{1}}, \mu_{1} \rangle, \ldots, \langle n^{k_{l}}, \mu_{l} \rangle) = \text{Proj}_{0}(\mathcal{X}, \mathcal{O}, \mathcal{L})(n^{k_{1}}, \mu_{1}) \circ \ldots \circ \text{Proj}_{0}(\mathcal{X}, \mathcal{O}, \mathcal{L})(n^{k_{l}}, \mu_{l}).
\]

\( \text{Proj} \) takes a state from a state trajectory and returns either one pair or an empty sequence \( \varepsilon \). If \( n^{k} \) is an occurrence of interest, then it returns \( \langle (n^{k}, \mu \upharpoonright \mathcal{X}) \rangle \). This means that, at \( n \), we consider exact values of variables in \( \mathcal{X} \). If \( n^{k} \) is not an occurrence of interest but, due to a KL form, the projection has to keep trace of a set \( \mathcal{L} \) of executed statements (even if the variables in that point are not of interest), then \( \text{Proj} \) returns \( \langle (n^{k}, \mu \upharpoonright \emptyset) \rangle \), meaning that we require the execution of \( n \), but we are not interested in the values of variables in \( \mathcal{X} \).

Trajectory projection allows us to define all the semantic equivalence relations characterizing on what a program and its slices have to agree due to the chosen criterion. Given two programs \( P \) and \( Q \), we can say that \( Q \) is a slice of \( P \) if it contains a subset of the original statements and \( Q \) is equivalent to \( P \) with respect to the semantic equivalence relation induced by the chosen slicing criterion.

**Definition A.4 (Binkley et al. [2006a]).** Let \( P \) and \( Q \) be executable programs and \( C = (\mathcal{I}, \mathcal{X}, \mathcal{O}, \psi) \) be a slicing criterion. Let \( \mathbb{L}_{P} \subseteq \mathbb{L} \) be the set of program points of \( P \) and
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Fig. 24. Execution trace of \( P \) in Example A.5 and of the static slice and of their projection.

\[ \mathcal{L} = \mathbb{L}_P \cap \mathbb{L}_Q \text{ if } \psi = \text{true}^{16} \quad (\mathcal{L} = \emptyset \text{ if } \psi = \text{false}). \] Then \( P \) is equivalent to \( Q \) w.r.t. \( C \) if and only if

\[ \forall \mu \in \mathcal{I}. \text{Proj}_{X,O,L}(\tau^\mu_P) = \text{Proj}_{X,O,L}(\tau^\mu_Q). \]

The function \( \mathcal{E} \) maps any criterion \( C \) to the corresponding semantic equivalence relation; hence, in this case, we write \( \langle P, Q \rangle \in \mathcal{E}(C) \).

**Example A.5.** Consider the Program \( P \) in the left of Figure 7; let the input for \( n \) be 2. Suppose we want to compute a non-iteration-count KL form of dynamic slicing, that is, \( C = (\{n \leftarrow \{2\}\}, \{i, s\}, (8, N), \text{true}) \). Namely, the variables of interest are \( i \) and \( s \), which are observed at the program point 8 each time it is reached, and the slice does not have to execute statements not executed in the original program. The program on the right of Figure 7 is a slice w.r.t. \( C \). In Figure 24, we have the execution trajectory of the original program (on the top), the execution trajectory of the candidate slice (in the middle), and the (same) projection of the two trajectories due to the chosen criterion (on the bottom).

The formal framework proposed in Binkley et al. [2006a, 2006b] represents different forms of slicing by means of a \((\sqsubseteq, \mathcal{E})\) pair: a syntactic preorder and a function from slicing criteria to semantic equivalences. The preorder fixes a syntactic relation between the program and its slices. In traditional slicing, slices are obtained from the original program by removing zero or more statements. This preorder is called traditional syntactic ordering, simply denoted by \( \sqsubseteq \), and it is defined as follows: \( Q \sqsubseteq P \iff \mathbb{L}_Q \subseteq \mathbb{L}_P \).

The second component \( \mathcal{E} \) fixes the semantic constraints that a subprogram has to respect in order to be a slice of the original program. As we have seen before, the equivalence relation is uniquely determined by the chosen slicing criterion determining also a specific form of slicing. This way, Binkley et al. are able to characterize eight forms of non-SIM slicing and 12 forms of SIM slicing.

Finally, this framework is used to formally compare the different notions of slicing. First of all, it defines a binary relation on slicing criteria \( \rightarrow \) [Binkley et al. 2006a]: Let

---

\[^{16}\text{Note that, when } \mathbb{L}_Q \subseteq \mathbb{L}_P, \text{ as we suppose in this article, then } \mathcal{L} = \mathbb{L}_Q. \text{ We provide the general definition since the original definition of dynamic slicing [Korel and Laski 1988] does not require that all the lines of } Q \text{ be included in } \mathcal{L}; \text{ however, our choice follows the paths taken in the original framework [Binkley et al. 2006b].} \]
Fig. 25. Given two forms \( A \) and \( B \), both (non-)\( \text{SIM} \), \( A \) is weaker than \( B \) if \( A \) is connected to \( B \) by a solid line and it is below \( B \). If \( A \) is non-\( \text{SIM} \) and \( B \) is \( \text{SIM} \), then \( A \) is weaker than \( B \) if \( A \) is connected to \( B \) by a dotted line and it is to the left of \( B \).

\[
C^1 = \langle I^1, X^1, O^1, \psi^1 \rangle \quad \text{and} \quad C^2 = \langle I^2, X^2, O^2, \psi^2 \rangle:
\]

\[
C^1 \rightarrow C^2 \quad \text{iff} \quad I^1 \subseteq I^2, X^1 \subseteq X^2, O^1 \subseteq O^2, \psi^1 \Rightarrow \psi^2. \tag{2}
\]

At this point, we say that a form of slicing \( (\subseteq, \mathcal{E}^1) \) is weaker than \( (\subseteq, \mathcal{E}^2) \) w.r.t. \( \rightarrow \) if and only if \( \forall C^1, C^2, \) slicing criteria such that \( C^1 \rightarrow C^2, \) and \( \forall P, Q, \) if \( Q \) is a slice of \( P \) w.r.t. \( (\subseteq, \mathcal{E}^2(C^2)) \), then \( Q \) is a slice of \( P \) w.r.t. \( (\subseteq, \mathcal{E}^1(C^1)) \) as well. In this case, we say that \( (\subseteq, \mathcal{E}^1) \) subsumes \( (\subseteq, \mathcal{E}^2) \). Following this definition, Binkley et al. show that all forms of slicing introduced in Binkley et al. [2006a] are comparable in the way shown in Figure 25, where the symbols \( S, C, D, SS, SC, \) and \( SD \) represent static, conditioned, dynamic, static \( \text{SIM} \), conditioned \( \text{SIM} \), and dynamic \( \text{SIM} \) types of slicing, respectively. Subscripts \( i, KL, \) and \( KL_i \) represent IC, KL, and KL\(_i\) forms of slicing, respectively; the absence of subscripts denotes the standard forms of slicing. In Figure 25, we explicitly provide the hierarchy concerning both \( \text{SIM} \) and non-\( \text{SIM} \) forms of conditioned slicing constructed in Binkley et al. [2006a].

ACKNOWLEDGMENTS

This work was funded partially by the EU project FP7-ICT-610582 ENVISAGE: Engineering Virtualized Services (http://www.envisage-project.eu), by the Spanish MINECO projects TIN2012-38137 and TIN2015-69175-C4-2-R, and by the CM project S2013/ICE-3006.

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Received December 2015; revised September 2016; accepted December 2016