Conservative fixpoint functions on a graph

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Conservative Fixpoint Functions
on a Graph

by

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Abstract

In this paper we present a derivation of a general solution for a class of programming problems. In these problems a function over the vertices of a directed graph is to be computed, being defined as a least fixed point of some monotonic operator. If this operator satisfies a certain restriction with respect to its image for a differential change in its argument, it is called conservative, and an elegant general solution may be derived. It is stipulated that a strictly calculational derivation is only possible if the level of abstraction is sufficiently high. To that end a modest extension to the functional calculus is proposed, including partial functions, and a few simple high level programming constructs are introduced. The program scheme obtained is applied to a particular example, for which so far no derivation, other than informal ones, is known to exist. The solutions presented are not new, but the calculational, abstract and compact technique of deriving them is meant to improve and complement the current techniques [KNU77, REY81, REM84]. It is believed to simplify the derivations for a wider algorithm class [EIJ92] than the one treated here.
1. Introduction

From a fundamental point of view solving fixed-point equations for functions on a graph is interesting, because they represent generalizations of recursive equations "par excellence": each function value may depend on an arbitrary, possibly empty, set of other function values. These predecessor or successor sets simply define the graph. Obviously, the cyclicity of recursive equations is then reflected by the cyclicity of the corresponding graph. On the other hand, these fixed-point equations are of great practical use. They arise in many different fields: topography, digital networks, electronic circuits, compiler theory, computer graphics, and process scheduling. Recently the importance of solving fixed-point equations on graphs was recognized by Cai and Paige [CAI89], who devoted a comprehensive paper with many examples to the subject, although their objectives differ from ours, e.g. automatic code generation.

The technique used in this paper to solve graph programming problems is based on three corner-stones: heuristics, abstraction and calculus. Some twenty-five years ago, E.W. Dijkstra was one of the first to advocate the use of proper heuristics [DIJ76]. For a more up-to-date treatment the reader is referred to [GAS88]. These heuristic principles, and some others, are applied wherever appropriate in this paper, and an attempt is made to indicate at crucial points in the derivation why certain design decisions were taken. In [EIJ92] the reader may find explicitly formulated "rules of thumb", applying more in particular to set, graph or fixed-point problems.

The need for abstraction is widely accepted [REY81, CAI89], although in some modern textbooks [MEH84, COR90] graph algorithms are treated in terms of low level data structures. In my opinion, abstraction is the only way to reason about more complicated algorithms like graph algorithms, allowing the methods for solving small problems [DIJ88] to be applicable without requiring major changes. However, the way to abstract is subject to discussion. Since we are interested in computing functions the most logical conclusion seems to be to introduce abstract functions as data structures into the basic formalism, for which we choose the guarded command language (GCL) [DIJ88]. Because in repetitions we would like these functions to grow towards their final value we also introduce partial functions. Finally, we need abstract sets, in the spirit of [REM84]. The details of this data abstraction are summarized in section 2, together with a touch of control abstraction [LIS77] by means of some convenient programming structures, the most important one being a
for-statement that repeats a statement list for all elements of a fixed set. A straightforward proof rule is expressed by the precondition for theorem, and section 2 is concluded with the expression accumulation theorem, needed to make efficient use of the for-statement.

The third cornerstone of our technique, a calculational style, is widely believed to be a prerogative of transformational programming, be it Imperative [MOR90], recursive [BUR77], functional [BIR88] or even relational [ROL90]. One of the goals of this paper is to show that programming with conventional Hoare triples can also be done purely calculationally. The required set and function calculus rules are introduced in the beginning of section 2. The resulting style is, admittedly, hybrid, with program notation on the one hand, and derivations using predicates on the other. On some occasions, e.g. in the expression accumulation theorem or the function lift rules, the influence of functional programming can be distinguished, and used to advantage.

The problem class under consideration is defined in section 3. Many members have been treated in literature, from minimal distance problems, dating as far back as 1959 [DIJ59, TAR81], reachability problems [REY81, CAI89], and least ancestor problem [REM86], to the capacity or tunnel problem [CAR79, REM85]. The derivation presented in section 4 solves the general problem. Compared to [CAI89, REM85], it has the advantage of being purely calculational, requiring inventiveness only at certain crucial points.

Many of the problems in the class can be treated very elegantly using fixed-point theory and a regular algebra of matrices [BAC75, CAR79, TAR81]. However, the present class of problems is more general, because these matrices carry less information than the fixed-point operators on functions in section 3. Therefore, in section 5, the general problem is instantiated to an example that, as far as I know, cannot be described by such a regular algebra: the problem of ascending reachability [REM85]. It does, however, turn out to be an instance of the problem class in [KNU77], but there a non-calculational proof of the solution algorithm is only given afterwards.

2. The SAL formalism

In the program notation of SAL (Set Algorithmic Language), which is an extension of GCL, all variables are assigned a - constant - type at the beginning of the block with respect to which they are local, by a declaration
of the form \( v : T \) (\( v \) a variable and \( T \) a type). Types are considered sets, supplemented with a collection of operators, satisfying certain laws.

In addition to the simple types \( \mathbb{Z} \) (integers) and \( \mathbb{B} \) (booleans), adapted from GCL, with the usual operators and laws, a separate, unstructured type \( U \) is introduced, the elements of which represent graph vertices. Apart from the array type constructor in GCL, we also have \( \mathcal{P} \) (powerset type constructor), supporting the usual set operators and laws, \( \to \) (function type constructor), and \( \to^* \) (partial function type constructor).

By convention sets are denoted with capitals, and \( u, v, w, x, y, \) and \( z \) are understood to be vertices. Unless stated otherwise, they are of standard type \( \forall v : \mathbb{P} \cdot U \). Lower case letters \( d, e, f, g, \) and \( h \) denote partial or total function variables. On the other hand, for predicates and other mappings \textit{sculpt} letters are used. Simple and generic types, and the powerset type constructor, are in \textbf{DOUBLE CAPITALS}. Finally, statement list names are in \textit{italics}.

Although in principle the Eindhoven Quantifier Notation (EQN) \cite{DIJ88} is adequate to formulate and prove the usual mathematical statements, I agree with \cite{BAC89} that for complicated problems it tends to get very lengthy and cumbersome. So whenever possible, the dummies are abstracted from, and calculations are carried out on the function or set level. If not, we use the quantifiers \( \forall \) (universal), \( \exists \) (existential), \( \downarrow \) (minimum) and \( \uparrow \) (maximum), in an EQN-like notation, with an explicit dummy, a domain and a term, all separated by colons. The same convention is used for sets, so \( \{ v : D \cdot v : \mathcal{E} \cdot v \} \) denotes the set of values \( \mathcal{E} \cdot v \) for those \( v \) that satisfy predicate \( D \). In order to avoid confusion between set notation on the one hand, and the notation of hints in proofs and assertions in algorithms on the other, we adopt the convention of replacing the hint and assertion brackets (\{\}) by the special brackets \{\}.

Total function \( f : F \rightarrow T \) has domain \( F \) and range \( T \), and the so-called domain operator \( d \) satisfies \( d \cdot f = F \). On the other hand, for a partial function \( g : G \rightarrow T \), \( d \cdot g \) denotes the subset of domain \( G \) on which \( g \) is defined. So partial functions can be considered total functions with a variable domain; they are needed as program variables. All function notations and properties in the remainder of this section apply to total as well as partial functions, and \( \to \) may denote either type constructor; \( f : F \rightarrow T \) and \( g : G \rightarrow T \) are supposed to be arbitrary functions. By convention, the domain of a named function is denoted by the corresponding name in \textbf{CAPITALS}, hence

\[
F = d \cdot f \quad \text{and} \quad G = d \cdot g.
\]
Instances of functions can be denoted as sets of pairs, and $\emptyset$ is the **empty function**. The constant function with domain $S$ and value $c$ for all $s \in S$ is denoted by $(S; c)$, and if $S = \{s\}$ is a singleton, this will be abbreviated to $(s; c)$, the **singleton function**. A constant function with standard domain $V$ will be written $c$, or $S^{-}$ if $S$ is an expression. The **lambda quantifier** offers a general notation for functions:

$$(\lambda i: D \cdot i: E \cdot i)$$

is the function which, for all values of dummy $i$ that satisfy $D \cdot i$, takes the value $E \cdot i$. In this notation a domain restricted function is defined by

$$(A \cdot f) = (\lambda i: i \in A \land F: f \cdot i), \quad (0)$$

for arbitrary set $A$.

The operators on functions are defined in terms of the operators on their respective ranges, lifted to the function level. Let $\star \in T \to T \to T$ be a binary infix operator. Then the lifted operator $\hat{\star}$ is defined by the **lift-meet rule**:

$$f \hat{\star} g = (\lambda i: i \in F \land G: f \cdot i \hat{\star} g \cdot i) \in F \land G \to T. \quad (1)$$

So far, rules (0) and (1) only allow us to restrain function domains. In programming, if a function is to be computed, we need to be able to extend a function domain, or to combine functions with disjoint domains. To that end we introduce the so-called **lift-split rule**, which, assuming $F \land G = \emptyset$, is given by

$$f \hat{\cdot} g = (\lambda i: i \in F \land U: \text{if } i \in F \cdot i \hat{\cdot} (\neg F \cdot g) \cdot \text{if } i \in G \cdot g \cdot i \hat{\cdot} f) \in F \land U \to T. \quad (2)$$

As it happens, the functions encountered in programming practice rarely have disjoint domains. In fact, a combination of (1) and (2) turns out to be very fruitful. An obvious definition for yet another lifted operator $\check{\star}$ would be

$$f \check{\star} g = (f \hat{\star} g) \cdot (\neg G \cdot f) \cdot (\neg F \cdot g)$$

$$= (\lambda i: i \in F \land U: \text{if } i \in F \land G \to f \cdot i \hat{\cdot} g \cdot i) \in F \land U \to T. \quad (3)$$
This rule is called the lift-join rule. Its power is that it generalizes both (1) and (2): using (0) and (1) we have

\[ f \circ g = (F \circ G \cdot f \circ g), \tag{4} \]

and if \( F \circ G = \emptyset \) (3) reduces to (2); in that case the operator \( \circ \) on the left-hand side may be chosen arbitrarily, or omitted if it is irrelevant.

As for operator properties, it turns out that commutativity, associativity, idempotency and distributivity simply carry over to the lifted operators. For commutative and associative operators \( \circ \) we sometimes need a generalized version of (3), which reads

\[ (\circ k, D \cdot k; h \cdot k) = (\lambda i : i \in (U k : D \cdot k; H \cdot k) : (\circ k : D \cdot k \land i \in H \cdot k; h \cdot k \cdot i)), \tag{5} \]

where \( h \cdot k \in H \cdot k \to T \) is a \( k \)-indexed family of functions.

Some function calculus rules, relating (0), (2) and (3), are summarized below. For arbitrary sets \( A \) and \( B \) we have the chaining rule

\[ (A \land B \cdot f) = (A \cdot (B \cdot f)), \tag{6} \]

If, in addition, \( \circ \in T \to T \to T \) is an arbitrary binary infix operator on the common range \( T \) of \( f \) and \( g \), the term split rule expresses

\[ (A \cdot f \circ g) = (A \cdot f) \circ (A \cdot g). \tag{7} \]

For idempotent operators \( \parallel \in T \to T \to T \) the domain disjunction rule

\[ (A \lor B \cdot f) = (A \cdot f) \lor (B \cdot f) \tag{8} \]

holds. If \( A \land B = \emptyset \), rule (8) is called domain split rule, and operator \( \parallel \) may be omitted. Rules (6), (7) and (8) are easily derived from (0), (2) and (3). Detailed proofs are given in [EIJ92]. An application of (5), the domain shift to term rule, reads

\[ (\circ k, i : D \cdot k \land i \in W \cdot k : ((i) \cdot h \cdot k)) = (\circ k : D \cdot k : (W \cdot k \cdot h \cdot k)), \tag{9} \]

with \( \circ \) commutative and associative, and \( W \cdot k \) \( k \)-indexed sets. It follows from
Finally, we introduce a shorthand for a quantification over a function term, comparable to the reduce (/) from the Bird-Meertens formalism [BIR88], in the same way as function composition (\( \odot \)) is comparable to the map (\( \ast \)). For an associative and commutative operator \( \cdot \) we define

\[
\theta \cdot f = (\theta i: i \in F: f \cdot i).
\]

(10)

We conclude with two domain split rules for \( \theta \cdot f \), with a straightforward proof:

\[
\theta \cdot (f \cdot g) = \theta \cdot f \odot \theta \cdot g, \quad \text{and}
\]

(11)

\[
\theta \cdot (f \odot g) = \theta \cdot f \odot \theta \cdot g.
\]

(12)

Next we turn from data abstraction to control abstraction [LIS77]. Firstly, a convenient miracle statement is introduced. To that end we define statement

\[
\text{out: } \text{Post} \cdot \text{in} \cdot \text{out}
\]

with meaning: variable out is assigned a value such that \( \text{Post} \cdot \text{in} \cdot \text{out} \) is validated. A detailed definition of this programming primitive is beyond the scope of this paper; the reader is referred to [MOR90].

A related programming primitive is element selection. Statement \( x: \in X \) is defined using the weakest precondition \( \text{wp} \) for arbitrary predicate \( P \), by

\[
\text{wp} \cdot (x: \in X) \cdot P \equiv X \neq \emptyset \land (\forall v: v \in X: P^X_v).
\]

For reasons of convenience two new structured statements are introduced...
into SAL, in addition to concatenation, if ... fi and do ... od. The as statement is an alternative statement comparable to the if statement, but if none of the guards is satisfied a skip is executed, instead of an abort. Apart from being an obvious counterpart of the if, it has the advantage of making programs shorter and more transparent, at least for the majority of graph algorithms considered. For a single guard it is defined by

\[\text{as } B \rightarrow SL \text{ sa is equivalent to } \text{if } B \rightarrow SL \ 0 \neg B \rightarrow \text{skip fi},\]

where \(B\) is a boolean expression and \(SL\) a statement list. Generalizations to more guards are obvious, and left to the reader. A proof rule in terms of Hoare triples without reference to the if statement is given by:

\[\begin{align*}
\{P\} \text{ as } B \rightarrow SL \text{ sa } \{Q\} & \text{ is valid if and only if } \\
\{P \land \neg B\} \ SL \ \{Q\} & \text{ and } P \land \neg B \Rightarrow Q \ \text{ hold.}
\end{align*}\]

Another structured statement to be introduced in SAL is the for statement, which is a repetition, where a variable takes the value of all elements of a given set, in arbitrary order. It is defined by

\[\begin{align*}
\text{for } u \epsilon W \rightarrow SL \ \text{ref} & \text{ is equivalent to } \\
& \begin{align*}
& [\ A: P \cdot W; \ u: W | A:=W; \ \ do & A \neq \emptyset \rightarrow \\
& u:=A; \ A:=A\setminus\{u\}; \ SL \ \ od \\
& ].
\end{align*}
\]

Evidently, the for statement not only turns out to be shorter than its GCL counterpart, but it also saves the introduction of the fresh set variable \(A\). The occurrence of \(u\) in the for clause for \(u \epsilon W\) is supposed to implicitly open a block enclosing the repetition in which \(u\) is declared as a local variable. Of course \(u\) may be used in statement list \(SL\), whereas \(A\) may not.

Application of the for-statement evidently offers no special advantages, if it replaces the corresponding repetition only after the latter is derived in full. Preferably, we would like to apply a proof rule directly to the for-statement itself. The following theorem offers one that is directly related to (14).
Precondition for-theorem

Let $W \in \mathcal{P}\cdot V$ be a program expression yielding a finite set, and let $u \in W$ and $A \in \mathcal{P}\cdot W$ be fresh variables. If there exists a predicate $P$ depending on $A$, but not on $u$, and a statement list $SL$ that does not refer to $A$, then

\[
\{u \in A \land P^A_{A \cup \{u\}}\} \land SL \land \{P\},
\]

implies the validity of

\[
\{P^A_W\} \lor \{u \in W \rightarrow SL \lor \forall \} \land \{P^A\}.
\]

This theorem is proved in appendix A. It is generally applicable, but a drawback is the use of predicates, making application still rather elaborate. On the other hand, the theorem below addresses the special case of refining a complex assignment statement directly to a for-statement. It turns out to be adequate and practical for many applications. Not surprisingly, it reminds us of the so-called reduction from functional programming [BIR88]:

Expression accumulation theorem

Let $r \in \mathbb{T}$ be a program variable, $W \in \mathcal{P}\cdot V$ a program expression yielding a finite set, $\mathcal{S} \in \mathcal{P}\cdot V \rightarrow \mathbb{T}$ and $\mathcal{I} \in V \rightarrow \mathbb{T}$ mappings, and let $\mathcal{G} \in \mathbb{T} \rightarrow \mathbb{T}$ be a binary infix operator, for some type $\mathbb{T}$. If for all $A \in \mathcal{P}\cdot W$, $\mathcal{F} \in \mathbb{T}$, and $u \in W \setminus A$

\[
\mathcal{S}(A \cup \{u\}) \cdot t = \mathcal{S}(A \cdot t) \odot \mathcal{I} \cdot u
\]

holds, then statement $r := \mathcal{S} \cdot W \cdot r$, with precondition $r = \mathcal{S} \cdot \mathcal{G} \cdot r$, is refined by

\[
\forall \ u \in W \rightarrow r := r \odot \mathcal{I} \cdot u \lor \forall \phi . \]

A proof outline is given in appendix A. For a detailed proof the reader is referred to [EIJ92]. Note that, contrary to the so-called accumulation function $\mathcal{S}$, it is essential that $\mathcal{I}$ does not depend on $A$ or $t$.

3. Specification of the conservative fixpoint function problem

Consider a directed graph $G = \langle V, S \rangle$, with vertex set $V \in \mathcal{P} \cdot W$, and successor function $S \in V \rightarrow \mathcal{P} \cdot V$, mapping each vertex to the set of its direct successor.
vertices. Alternatively, the predecessor function $P \in V \rightarrow P \cdot V$ may be used instead of $S$, the relation between the two being given by

$$(\forall u,v:: v \in S \cdot u \equiv u \in P \cdot v).$$

Often $S$ and $P$ are generalized to functions of type $P \cdot V \rightarrow P \cdot V$, according to

$$S \cdot A = (u: v \in A: S \cdot v),$$

and analogously for $P$. Denoting both the standard and generalized versions of $S$ and $P$ with the same symbol rarely causes confusion. Both $S$ and $P$ are strict with respect to $\emptyset$, monotonic, and uni-disjunctive (universally disjunctive).

Let $(\mathbb{T}, \preceq)$ be a totally ordered, complete lattice, with bottom 1, supremum $\mathbb{U}$ and infimum $\mathbb{N}$ [DAV90]. We consider partial functions $V \rightarrow \mathbb{T}$ on the vertices of the graph. Then, using lift-join rule (3), we define lifted versions $\mathbb{U}$ and $\mathbb{N}$ of supremum and infimum, respectively. These, in turn, define in a natural way a partial order on $V \rightarrow \mathbb{T}$ (in fact a complete lattice [DAV90]), consequently denoted by $\preceq$, though itself not being defined by (3), but by

$$d \preceq e \equiv DSE \wedge \left(\forall v: v \in D: d \cdot v \preceq e \cdot v\right).$$

For this partial order $(A, \preceq)$ and $d \preceq e$ turn out to be monotonic.

We are asked to determine the least fixed point $f \in V \rightarrow \mathbb{T}$ of some operator $\mathcal{F} \in (V \rightarrow \mathbb{T}) \rightarrow (V \rightarrow \mathbb{T})$. So $f$ satisfies

$$\mathcal{F} \cdot f = f,$$

$$(\forall e: \mathcal{F} \cdot e = e: f \preceq e).$$

Operator $\mathcal{F}$ is assumed to be monotonic, hence it satisfies

$$d \preceq e \Rightarrow \mathcal{F} \cdot d \preceq \mathcal{F} \cdot e.$$  

Often recursive formulations can be derived from explicit expressions, using fixed-point theorems like that of Knaster-Tarski [DAV90]. In section 5 an instance of (20)-(21) is derived using the theorem below.

**Generalized fixed-point theorem**

Let $\mathcal{F} \in L \rightarrow L$ be a uni-disjunctive operator on the complete lattice $(L, \preceq, \mathbb{U}, \mathbb{N})$, and let $e$ be an element of $L$. If $r$ is given by
then \( r \) is uniquely defined as the least (pre)fixed point of \((e \cup) \cdot \hat{a}\), i.e.

\[
e \cup \hat{a} \cdot r = r, \quad \text{and} \quad (\forall s: s \in (e \cup) \cdot \hat{a} \cdot s \Rightarrow r \in s).
\]

Here \((e \cup) \cdot \hat{a}\) corresponds to \( S \) in (20)-(21). A proof is given in appendix B.

For the special class of problems considered here, \( S \) also satisfies a restriction on its image, with respect to a "differential" change in its argument, a kind of "limited growth"-property, expressed by

\[
S \cdot (e \cdot (x: c)) \subseteq S \cdot \hat{c}, \quad \text{for } x \in E.
\]

An operator satisfying (26) is called conservative, and so are its fixpoints.

It is not hard to prove that, if one restricts the context-free grammars from \([\text{KNU77}]\) to generate only graph problems, then fixed-point functions turn out to be conservative precisely if their generating functions are superior (see \([\text{KNU77}]\)). Though restricted to graphs, the present formalism is more general with respect to operators and function domains.

The class of programming problems defined by (18)-(26) can not necessarily be mapped to a regular algebra of matrices \([\text{BAC75, CAR79, TAR81}]\). On the other hand, if the regular algebra satisfies the property that the unit of its multiplication is top \((T)\) of the lattice \((T, E)\), a property that corresponds to (26), it is an instance of the problem class defined above.

The required \( f \) computation is now accomplished by realizing postcondition

\[\mathcal{R}: \quad g = f.\]

where \( g \in V \Rightarrow T \) is a partial function variable.

4. Derivation of an algorithm

Because the quantity to be computed is a least fixed point we approximate it from below (cf. linear search) and propose the simplest invariant possible:
\( P_0: \ g \notin f. \)

From \( P_0 \) immediately follows the first approximation to the desired algorithm, using the miracle statement:

\[
g := \mathcal{I} \;
\text{do}\ g \in f \rightarrow
\text{do} g \in d \notin f; \ g := d
\text{od},
\]

which is obviously (partially) correct. However, in this algorithm \( f \) is an unknown quantity, which should be eliminated. Again it has to be approximated from below, according to the derivation

\[
f = \{ (20) \} \mathcal{F} \cdot f \nsimeq \{ P_0; (22) \} = \mathcal{F} \cdot g, \quad (27)
\]

so it appears useful to introduce variable \( h \in \mathcal{V} \rightarrow \mathcal{T} \), specified by

\( P_1: \ h = \mathcal{F} \cdot g. \)

Naturally we want to replace \( f \) everywhere in the algorithm by \( h \). Since, from (27), this means strengthening the guard and selection statement, this has no consequences for the invariance of \( P_0 \), so we obtain as second approximation:

\[
g, h := \mathcal{F}, \mathcal{F} \cdot \mathcal{G};
\text{do}\ g \in h \rightarrow
\text{do} g \in d \in h; \ h := \mathcal{F} \cdot d; \ g := d
\text{od}.
\]

However, strengthening the guard implies weakening the postcondition, so we ought to make sure that \( \mathcal{R} \) still holds upon termination. Negation of the guard could mean that \( g \) and \( h \) are not comparable, so in order to still conclude something useful we are forced to introduce invariant

\( P_2: \ g \notin h. \)

Then the postcondition remains valid, for we have
\begin{align*}
g \not\in h & \\
\equiv \{ \mathcal{P}_2; \mathcal{P}_1 \} & \\
g = \mathcal{F} \cdot g & \\
\Rightarrow \{ (21) \} & \\
f \not\in g & \\
\equiv \{ \mathcal{P}_0 \} & \\
\mathcal{R}. & \quad (28)
\end{align*}

Evidently, \( \mathcal{P}_0 \) and \( \mathcal{P}_1 \) remain valid in the latest algorithm. The validity of \( \mathcal{P}_2^{g,h} \) just after the selection follows from \( \mathcal{P}_1 \) and \( \mathcal{F} \) monotonicity (22):

\begin{align*}
d \not\in h = \mathcal{F} \cdot g \not\in \mathcal{F} \cdot d,
\end{align*}

so the abstract, second approximation turns out to be (partially) correct.

So far everything has gone smoothly, requiring hardly any inventiveness. But at this point we take an important design decision. The only drawback of the second approximation seems to be the lack of clarity about the efficiency, or for that matter, termination at all. We decide to settle this question rigorously, and strive for a solution where domain \( G \) of \( g \) is explicitly extended with one vertex \( x \in H \), hence \( x \notin G \), in each turn of the repetition, and choose for \( d \cdot x \) the greatest value permitted by upper bound \( h \) for \( d \), i.e. \( h \cdot x \). Such a refinement would consequently result in the third approximation:

\begin{verbatim}
g,h:= \emptyset, \mathcal{F} \cdot \emptyset; do H \setminus G \neq \emptyset \rightarrow 
   x:= H \setminus G; c:= h \cdot x; 
   d:= (x \cdot c); h:= \mathcal{F} \cdot d; g:= d od.
\end{verbatim}

Note that domains \( G \) and \( H \) of functions \( g \) and \( h \), respectively, need not be adjusted separately. Variable \( c \in \mathcal{F} \) is only introduced as a shorthand. However, the above algorithm is only a true refinement of the previous one if

\( \mathcal{P}_3: \ (G \uplus h) = g \)

holds, for this implies that \( g \not\in h \) is equivalent to \( H \setminus G \neq \emptyset \), using (19).
Invariant \( P3 \) means that \( h \) equals \( g \) on domain \( G \) of the latter, see (0).

Can we keep \( P3 \) invariant? It holds initially, but before checking its invariance, we first mention a convenient property of \( d \), directly following from conservativity property (26) and \( P1 \):

\[
\mathcal{F} \cdot d \notin h \cup \mathcal{G}.
\]  

This also implies \( c = \mathcal{F} \cdot g \cdot x \notin \mathcal{F} \cdot d \cdot x \subseteq h \cdot x \cup c = c \), proving

\[
\mathcal{F} \cdot d \cdot x = c.
\]  

For the invariance of \( P3 \) we then derive

\[
P3 \rightarrow \mathcal{F} \cdot d.
\]  

\[
\mathcal{F} \cdot d \iff h \cup \mathcal{G}.
\]  

In other words: all \( h \cdot y \) for \( y \in G \) (\( \sim G \) is \( G \)'s complement with respect to \( V \)) should be at most all \( g \cdot z \) values. Establishing in turn the invariance of \( P4 \):

\[
U \cdot (\sim G \cdot h) \subseteq \Pi \cdot g.
\]

where

\[
P4: U \cdot (\sim G \cdot h) \subseteq \Pi \cdot g.
\]

In other words: all \( h \cdot y \) for \( y \in G \) (\( \sim G \) is \( G \)'s complement with respect to \( V \)) should be at most all \( g \cdot z \) values. Establishing in turn the invariance of \( P4 \):

\[
U \cdot (\sim D \cdot F \cdot d) \subseteq \Pi \cdot d
\]

\[
\iff d = g \cdot (x; c); \text{reduce domain split (11); } \sim D \subseteq \sim G
\]
If follows that the invariance of $\mathcal{P}_4$ in turn is guaranteed if $x$ is chosen such that its $h$ value $c$ is maximal on domain $H \mathcal{G}$! Now a maximum search is a rather trivial matter; however, since a witness $x$ of the maximum is required it gives us the opportunity to demonstrate the use of the precondition for theorem.

If, in proof rule (15), we use for invariant $\mathcal{P}$

$$\mathcal{P}: \quad U \cdot (\neg \mathcal{G} \setminus A \leftarrow h) \in h \cdot x \land x \in H \mathcal{G},$$

then we conclude that for $A = \emptyset \mathcal{P}$ implies the last member of (33), whereas for $A = H \setminus \mathcal{G}$ the restricted domain function collapses, hence $x : \in H \mathcal{G}$ suffices as initialization. Next, we evaluate precondition $y \notin A \land \mathcal{P}_A^A(y)$ of the for-body statement list $SL$ to be determined, first under the condition $h \cdot y \in h \cdot x$:

$$y \notin A \land U \cdot (\neg \mathcal{G} \setminus (A \cup \{y\}) \leftarrow h) \in h \cdot x \land h \cdot y \in h \cdot x \land x \in H \mathcal{G}$$

which equals $\mathcal{P}$. Alternatively, under the condition $h \cdot x \in h \cdot y$ we derive

$$y \notin A \land U \cdot (\neg \mathcal{G} \setminus (A \cup \{y\}) \leftarrow h) \in h \cdot x \land h \cdot x \in h \cdot y \land x \in H \mathcal{G}$$

so it appears that if we perform $x := y$ in the second alternative, both yield $\mathcal{P}$.
In their postcondition. Hence, as proof rule (13) can be applied, and picking up the pieces we consequently find the fourth approximation

\[ g, h := \emptyset, \mathcal{F} \cdot \emptyset; \]
\[ \text{do } H \setminus G \neq \emptyset \rightarrow \]
\[ \quad x := H \setminus G; \]
\[ \quad \text{for } y \in H \setminus G \rightarrow \text{ as } h \cdot x \subseteq h \cdot y \rightarrow x := y \text{ sa rof}; \]
\[ \quad c := h \cdot x; \quad h := \mathcal{F} \cdot (g \cdot (x; c)); \quad g \cdot x := c \]
\[ \text{od,} \]

where we eliminated d and simplified the g assignment. This algorithm is the end of the line if we have no further information about \( \mathcal{F} \). Note that it is \( O(|\mathcal{F}|^2) \), if the assignment to h can be performed in \( O(|\mathcal{F}|) \) time.

In many problem instances \( \mathcal{F} \) satisfies split property

\[ \mathcal{F} \cdot (e \cdot (x; c)) = \mathcal{F} \cdot e \cdot (S \cdot x \cdot K \cdot c \cdot x), \quad (34) \]

where \( K \) is an indexed family of mappings, with \( K \cdot a \cdot u \in K \cdot a \cdot u \cdot \mathcal{T} \) for all \( a \in \mathcal{T} \) and \( u \in V \), and \( \ast \) is an associative binary infix operator. If \( \mathcal{F} \) does not satisfy (34), one could say there is something wrong with the particular choice of the graph, because the successor function is supposed to precisely reflect the dependency relations between function values in all vertices.

If (34) does hold, the h assignment becomes

\[ h := h \cdot \ast (S \cdot x \cdot K \cdot c \cdot x), \quad (35) \]

and this allows us to refine the assignment further, using the expression accumulation theorem with accumulation function \( \mathcal{F} \cdot A \cdot t = t \cdot \ast (A \cdot K \cdot c \cdot x) \). Hence we check property (16) of the theorem:

\[ \mathcal{F} \cdot (A U(y)) \cdot t \]
\[ = \langle \text{definition of } \mathcal{F} \rangle \]
\[ t \cdot \ast (A U(y)) \cdot X \cdot c \cdot x) \]
\[ = \langle \text{domain split of } \ast (8), y \notin A; \text{ choose } \ast \text{ as arbitrary operator} \rangle \]
\[ t \cdot \ast (A \cdot X \cdot c \cdot x) \cdot \ast ((y)) \cdot X \cdot c \cdot x) \]
\[ = \langle \text{definition of } \mathcal{F} \rangle \]
\[ \mathcal{F} \cdot A \cdot t \cdot \ast ((y)) \cdot X \cdot c \cdot x), \]
so \( f \cdot y \) is chosen to be \((y)\cdot K \cdot c \cdot x\), which is indeed independent of \( A \) and \( t \). Finally, \( G \cdot \theta \cdot t = t \) is a tautology, so the restriction on the precondition of (35) holds. Hence (35) is refined by the for-statement

\[
\text{for } y \in S \cdot x \to \ h := h \cdot ((y) \cdot K \cdot c \cdot x) \ \	ext{ref},
\]

which may be expanded using (3) and (13). Also note that, since upon termination of the outer repetition \( f = g = h \) holds, \( g \) has become superfluous, apart from its domain \( G \). Substituting these findings into the fourth approximation we arrive at the final, fifth approximation

\[
\begin{align*}
\begin{array}{l}
\text{for } y \in S \cdot x \to \ h := h \cdot ((y) \cdot K \cdot c \cdot x) \ \	ext{ref},
\end{array}
\end{align*}
\]

Here, as in many applications, it is possible to implement a partial function using an array of size \#V, with "blanks" indicating vertices outside its domain, making membership tests for function domains \((y \in H, y \in K \cdot c \cdot x)\) easy, \( O(1) \). Set \( H \setminus G \) - remember that \( G \setminus H \) - may be implemented with a left adjusted array (a "stack"), making element selection also \( O(1) \), provided that in the maximum \( x \) search each new \( x \) value after \( x := y \) is swapped to the back. Note that \( G \) itself is then obsolete. The details of the implementation phase are left to the reader, since the translation is a rather mechanical process. Also various optimizations, e.g. initializing \( h \) outside \( d \cdot (\emptyset \cdot \theta) \) with the left unit of \( \circ \) and simplifying the central as-statement, take us beyond the scope of the paper.

Correcting for various differences in assumptions and notation the algorithm derived above corresponds to those in [BAC75, CAR79, TAR81, KNU77].
5. An example: ascending reachability

At this point we consider an application of the solution scheme derived in
the previous section, called ascending reachability. This problem was proposed
and solved in [REM85]. The problem is to determine the set of vertices
reachable from a given set \( B \) via an ascending path. In this context "ascending
path" means a path with the edge labels in ascending order, the edges being
labelled by means of a weight function \( t \in V \times V \rightarrow \mathbb{Z} \). On graph edges \( t \) is
assumed to be finite.

It would not be fair to formally define the problem in a way most suited to
formal manipulation, as one is tempted to do. Transformation of the "most
natural" problem formulation to a suitable form is, in my opinion, an
important part of the problem, and by no means the easiest. If possible, it
should be based on heuristic principles that are more widely applicable.

To express the set of ascending reachables formally in a natural way, we
need a path formalism. Paths are considered non-empty sequences of vertices,
where all pairs of successive vertices \( x,y \) in such a sequence satisfy \( y \in S \cdot x \).
Thus \([x]\) is the zero-edge ("empty") path starting and ending in \( x \), \([x,y]\) is
the one-edge path from \( x \) to \( y \), and so on. In the following \( p \) is understood to
be a path in the graph. The vertices on path \( p \) are numbered, from 0 up to and
including \( \#p \), the path length, and \( p \cdot i \), with \( 0 < i < \#p \), is the \( i \)-th vertex. If
the last vertex of a path coincides with the first of another path, these
paths may be concatenated with operator \( \# \). Note that this is not the usual
sequence concatenation! This brief summary will do for our purposes.

In terms of the path formalism the set of ascending reachables is given by

\[
R = \{ p : p \cdot 0 \in B \land aac \cdot p : p \cdot \#p \},
\]

with predicate \( aac \) being defined recursively by

\[
\begin{align*}
aac \cdot [u] ,

\text{(38)} \\
aac \cdot (p \# [p \cdot \#p , u]) & \equiv aac \cdot p \land t \cdot (p \cdot \#p) \cdot u \geq m \cdot p ,
\text{(39)}
\end{align*}
\]

for \( u \in S \cdot (p \cdot \#p) \), and function \( m \) in turn being defined by

\[
\begin{align*}
m \cdot [u] & = -\infty ,

\text{(40)} \\
m \cdot (p \# [p \cdot \#p , u]) & = t \cdot (p \cdot \#p) \cdot u ,
\text{(41)}
\end{align*}
\]
again for \( u \in S \cdot (p \cdot \#p) \), i.e. the weight of the last edge of path \( p \), and, if restricted to ascending paths, indeed the maximum label weight on the path.

If one wishes to introduce a regular algebra \([BAC75]\) to describe (38)-(41), a homomorphism is needed on the data structure of paths with concatenation, corresponding to the regular algebra product. However, (39) tends to violate associativity. In any case, \( t \cdot u \cdot v \) interpreted as a matrix is not adequate to model the problem, and at least tupled matrix elements (e.g. ascendingness, first edge label, last edge label) are required. In this way associativity can be restored, but then the regular algebra product lacks a proper unit, and it is hard to define a regular algebra sum operator, especially one over which the product should distribute. If there is a way out, it may be rather messy.

On the other hand, the present problem is easily modeled using Knuth's formalism \([KNU77]\). The interested reader may verify that "superior" functions

\[
g_{uv} \cdot c = \begin{cases} 
  t \cdot u \cdot v < c & \rightarrow 0 \\
  t \cdot u \cdot v > c & \rightarrow t \cdot u \cdot v
\end{cases}
\]

will do the trick. The superiority is expressed by the fact that \( g_{uv} \cdot c \succ c \).

Returning to our problem we note that definition (37) has two drawbacks. Firstly, it requires the complicated path formalism, and secondly, it does not have the shape of a recursive equation, which proved to be so fruitful in last section's derivation, and for other graph problems \([EIJ92]\). It turns out that a simple recursive equation of that kind does not even exist for \( R \): the ascending reachability of a given vertex \( x \) cannot be expressed in that of its predecessors. Intuitively, we need the "degree of ascending reachability" of \( x \)'s predecessors, in order to see whether paths to them can be extended with an edge to \( x \). Therefore we apply a technique called information extension \([EIJ92]\), and introduce a function \( f \in V \rightarrow \mathbb{Z} \), defined by

\[
f = (\lambda v: (\exists p:: ap \cdot p \cdot v): (\downarrow p: \downarrow ap \cdot p \cdot v: m \cdot p)), \tag{42}
\]

using the shorthand

\[
ap \cdot p \cdot v = p \cdot 0 \in B \land p \cdot \#p = v \land aoc \cdot p. \tag{43}
\]

The minimum quantor in the term of (42) carries information about the most favourable incoming path in \( v \), while the severe domain restriction serves two purposes: considering only interesting vertices with finite \( f \) values increases the computing efficiency, and in addition \( R \) is conveniently equal to \( d \cdot f (=F) \).
Function $f$ contains more information than $R$ in the sense that $R$ can be expressed in terms of $f$, whereas the reverse is not possible.

Next, we try to transform the, indeed cumbersome, expression for $f$ into a recursive fixed point equation, by first expressing it as a countably infinite supremum, so that subsequently a fixed point theorem can be applied, like the one in section 3. We derive

\[ f = \{ (42); \text{one-point rule} \} \]
\[ = (\lambda v: (\exists p,w:: ap\cdot p\cdot w \land v \in (w))\): (\exists p,w:: ap\cdot p\cdot w \land v \in (w)\): m\cdot p) \]
\[ = \{ \text{generalized lift-join rule (5), with } i:=v \text{ and } k:=(p,w) \} \]
\[ = (\exists p,w:: ap\cdot p\cdot w:: (w:: m\cdot p)) \]
\[ = \{ \text{domain split over all path lengths } k; \text{ introducing } Z \text{ below} \} \]
\[ (\exists k:: 0 \leq k:: Z\cdot k) , \quad (44) \]
\[ \text{with } \quad Z\cdot k = (\exists p,w:: ap\cdot p\cdot w \land \#p=k:: (w:: m\cdot p)) . \quad (45) \]

If (44) is to be an explicit fixed-point expression like (23), $Z\cdot k$ should take the form $d^k\cdot (Z\cdot 0)$, with $d$ a uni-disjunctive operator. Substituting $k=0$ into (45), we find $p=[w]$, evaluate $\alpha \alpha \cdot [w]$ and $m\cdot [w]$ using (38) and (40), use (5) with $k:=w$, $D\cdot k:= w\in B$, and $h\cdot k:=(w:: -\omega)$, and arrive at

\[ Z\cdot 0 = (\exists w:: ap\cdot [w]\cdot w:: (w:: m\cdot [w])) = (\exists w:: w\in B:: (w:: -\omega)) = (B:: -\omega) . \quad (46) \]

Next, trying to express $Z\cdot (k+1)$ in terms of $Z\cdot k$ we derive

\[ Z\cdot (k+1) \]
\[ = \{ (45) \} \]
\[ = (\exists p,w:: ap\cdot p\cdot w \land \#p=k+1:: (w:: m\cdot p)) \]
\[ \text{with } \quad Z\cdot (k+1) = (\exists p,w:: ap\cdot p\cdot w \land \#p=k+1:: (w:: m\cdot p)) . \quad (47) \]
\[ \{ \text{identifying domain and term of (45) using (5)} \} \]

\[ (\downarrow v, w: v \in d \cdot (Z \cdot k) \land t \cdot v \cdot w \supset Z \cdot k \cdot v \land w \in S \cdot v: (w: t \cdot v \cdot w)) \]

\[ \{ \text{introduce } L \text{ below; calculus} \} \]

\[ (\downarrow v, w: v \in d \cdot (Z \cdot k) \land w \in S \cdot v \land L \cdot (Z \cdot k) \cdot v: (w: t \cdot v \cdot w)) \]

\[ \{ \text{shift } w \text{ to the term (9); definition of } A \text{ below} \} \]

\[ A \cdot (Z \cdot k), \quad (47) \]

where

\[ L \cdot e \cdot v = \{ z: t \cdot v \cdot z \supset e \cdot v: z \}, \quad \text{and} \]

\[ A \cdot e = (\downarrow v: v \in E: (S \cdot v \land L \cdot e \cdot v \supset t \cdot v)). \quad (49) \]

From (44), (46) and (47) we conclude

\[ f = (\downarrow k: 0 \leq k: A^k \cdot (B: \models \omega)). \quad (50) \]

In order to prove that \( A \) is uni-disjunctive we use the last-but-two expression in derivation (47), abbreviating \( w \in S \cdot v \) to \( P \) and \( (w: t \cdot v \cdot w) \) to \( T \). We derive

\[ A \cdot (\downarrow i: D \cdot i: E \cdot i) \]

\[ \{ \text{definition of } A; \text{ domain of a generalized lift according to (5)} \} \]

\[ (\downarrow v, w: v \in (U i: D \cdot i: E \cdot i) \land t \cdot v \cdot w \supset (\downarrow i: D \cdot i: E \cdot i \cdot v \land P: T) \]

\[ \{ \text{property of } U; \text{ term of generalized lift according to (5)} \} \]

\[ (\downarrow v, w: (\exists i: D \cdot i: v \in E \cdot i) \land t \cdot v \cdot w \supset (\downarrow i: D \cdot i \land v \in E \cdot i: E \cdot i \cdot v) \land P: T) \]

\[ \{ \text{property of } \downarrow \} \]

\[ (\downarrow v, w: (\exists i: D \cdot i \land v \in E \cdot i: t \cdot v \cdot w \supset E \cdot i \cdot v) \land P: T) \]

\[ \{ \text{generalized domain disjunction} \} \]

\[ (\downarrow i: D \cdot i: (\downarrow v, w: v \in E \cdot i \land t \cdot v \cdot w \supset E \cdot i \cdot v \land P: T)) \]

\[ \{ \text{definition of } A \} \]

\[ (\downarrow i: D \cdot i: A \cdot (E \cdot i)). \]

Hence all conditions for application of the generalized fixed point theorem are satisfied, and \( f \) in (50) is of the required shape (23). It follows that \( f \) is the least fixed point of an operator \( \mathcal{F} \) defined by

\[ \mathcal{F} \cdot e = (B: \models \omega) \downarrow A \cdot e, \quad (51) \]

and as such matches (20)-(21), the starting point of the algorithm derivation.
The uni-disjunctivity of $\mathcal{A}$ implies monotonicity of $\mathcal{F}$, if we take $\downarrow$ for $\cup$, and $\triangleright$ for $\land$. It remains to check properties (26) and (35). We derive

$$\mathcal{F}.(e \cdot (x; c))$$

$$= \{ (51); (49) \}$$

$$= \{ (48); (49); (51) \}$$

$$= \{ \text{domain split; (6)} \}$$

$$= \{ \text{one-point rule; chaining \& (6)} \}$$

so we recognize the required pattern of property (35) if we take $\downarrow$ for $\cup$, and

$$K \cdot c \cdot x = (\mathcal{L} \cdot c \cdot x \uparrow t \cdot x),$$

and at the same time, realizing that $K \cdot c \cdot x \cdot y = t \cdot x \cdot y$ is defined only on domain $\mathcal{L} \cdot c \cdot x$ (i.e. for $t \cdot x \cdot y \triangleright c$), it follows that (26) also holds. It seems we can instantiate the solution algorithm from the previous section with the above, refine the $h$ initialization, from (46) and (51), by a $\text{for}$-statement, and after some minor simplifications we arrive at

$$\begin{align*}
\lbrakk & G : \text{P} \cdot V | G := \emptyset; \\
& \text{for } x \in B \rightarrow h \cdot x := \infty \text{ \& \& (51); (49); (6); (51)}; \\
& \text{do } H \setminus G \neq \emptyset \rightarrow \\
& \lbrakk \lbrakk & x : V; c : \mathcal{Z} | x := H \setminus G; \\
& \text{for } y \in H \setminus G \rightarrow \text{ as } h \cdot x \triangleright h \cdot y \rightarrow x := y \text{ \& \& (51); (49); (6); (51)}; \\
& c := h \cdot x; \\
& \text{for } y \in S \cdot x \rightarrow \\
& \text{ as } y \in H \land t \cdot x \cdot y \triangleright c \rightarrow h \cdot y := h \cdot y \uparrow t \cdot x \cdot y \land y \in H \land t \cdot x \cdot y \triangleright c \rightarrow h \cdot y := t \cdot x \cdot y \text{ \& \& (51); (49); (6); (51)}; \\
& G := G \cup (x); \\
& \rbrakk \rbrakk ; \\
& \text{od } \uparrow h = f \land G = R \rbrakk \\
\rbrakk.$$

After the repetition we have $G = H$, and in accordance with the remark at the beginning of this section the latter is equal to $R$, so we need not calculate $R$. 
separately after all. One more remark: if \( h \) is initialized to \( +\infty \) outside \( B \), the central \( \text{as} \)-statement reduces to

\[
\text{as } t \cdot x \cdot y \rightarrow c \rightarrow h \cdot y := h \cdot y \downarrow t \cdot x \cdot y \text{ sa,}
\]

and upon termination \( h \) equals \( f \), extended to the entire set \( V \).

Comparing our solution to the one in [REM85], we observe that they are identical, except for the fact that Rem represents the domain part of function \( h \) with a left adjusted array, in which the order of the arguments matches the order of the vertices in \( H \cdot G \). That does not save memory, however, the worst case size of \( H \) still being \( O(#V) \).

The proof in [REM85] turns out to be entirely non-calculational, in fact verbal, though based on a complete set of more or less formal invariants, and it does not have the shape of a derivation. A flaw in the reasoning is the claim that the weight of any ascending path to the "minimal" vertex \( x \) in \( H \cdot G \) must be at least \( (\downarrow u: u \in H \cdot G: h \cdot u) \), because any such path has an initial part in \( G \) ending in a vertex of \( H \cdot G \). This argument does not stick, because on \( H \cdot G \) \( h \) is not necessarily equal to \( f \), all these \( h \cdot v \) values may still decrease.

It should be noted that the optimized solution using heaps in [REM85] is incomplete. Apart from restoring the heap-organized set \( H \cdot G \) each time a vertex is added to \( H \), also statement \( h \cdot y := h \cdot y \downarrow t \cdot x \cdot y \) destroys the heap-order. Repairing this takes \( O(\log #V) \). It can be carried out many times for one and the same vertex, but is in total bounded by the number of edges \( #E \). So the entire solution then becomes \( O((#V + #E) \cdot \log #V) \). Surprisingly, this is no worse than Rem claims, but then again his - faulty - solution is actually \( O(#V \cdot \log #V + #E) \). As stated before, our solution is \( O(#V^2) \), which can be slightly worse than the heap solution, depending on the density of the graph.

6. Acknowledgements

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7. References


Appendix A: Proof of the precondition for-theorem

Use \( P \) as invariant of the corresponding do statement in (14), and \( \#A \) as variant function, then the correctness of (15) follows easily from the annotated version of the algorithm below.

\[
\begin{array}{l}
\left[ \begin{array}{l}
A: \ P \cdot W; \ u: W; \quad \{P_A^W\} \\
A:= W; \quad \{P\} \\
do \ A \not\in \emptyset \rightarrow \{A \not\in \emptyset \land P\} \\
u: \in A; \quad \{u \in A \land P\} \\
A:= A \setminus \{u\}; \quad \{u \in A \land P^A_{A \cup \{u\}}\} \\
SL \quad \{P\} \\
od \quad \{P \land A=\emptyset, \text{ hence } P^A_{\emptyset}\} \\
\end{array} \right. \\
\end{array}
\]

All transitions are trivial, only when applying the definition of \( u: \in A \) the independence of \( P \) on \( x \) is crucial, while the decrease of the variant function follows from \( \#(A-(u)) = \#A-1 \) (for \( u \in A \)), provided that SL does not refer to \( A \).

Proof outline for the expression accumulation theorem

It is easy to see that the desired refinement, say \( FOR \), of the assignment statement \( r:= \emptyset \cdot W \cdot r \) should satisfy
\( \{ r = \mathcal{G} \cdot \emptyset \cdot \rho \land W = \Omega \} \text{ FOR } \{ r = \mathcal{G} \cdot \Omega \cdot \rho \} \) \hspace{1cm} (A0)

for all values of the logical variables \( \rho \) and \( \Omega \). The validity of the precondition requirement has been taken into account. Hoare triple (A0) suggests a refinement of FOR by a repetition with invariant

\( \mathcal{G} : r = \mathcal{G} \cdot (\Omega \Delta) \cdot \rho, \) \hspace{1cm} (A1)

with a set variable \( A \), initially set to \( W \), shrinking to \( \emptyset \). The correctness of the repetition below follows from the annotation in brackets:

\[
\begin{align*}
\text{if } & A : W; \, u : W \mid \\
& A := W; \\
& \{ r = \mathcal{G} \cdot \emptyset \cdot \rho \land W = \Omega \} \\
& \text{do } A \neq \emptyset \rightarrow \\
& \quad \text{u} \in \text{A}; \\
& \quad \{ u \in \text{A} \land r = \mathcal{G} \cdot (\Omega \Delta) \cdot \rho \} \\
& \quad A := A \setminus \{ u \}; \\
& \quad \{ u \in \text{A} \land r = \mathcal{G} \cdot (\Omega \Delta (A \cup \{ u \})) \cdot \rho \} \\
& \quad r := r \odot \text{\| u} \\
& \quad \{ \text{u} \in \text{A} \quad \land r = \mathcal{G} \cdot (\Omega \Delta) \cdot \rho, \text{ so } r = \mathcal{G} \cdot \Omega \cdot \rho \} \\
\text{od}
\end{align*}
\]

The reader is invited to carefully examine all steps. Now, applying the definition of the for-statement (14), we find that (17) is indeed a correct refinement of FOR, and hence of \( r := \mathcal{G} \cdot W \cdot r \).

**Appendix B: Proof of the generalized fixed point theorem**

First we prove that \( r \), as defined in (23), satisfies (24). We derive

\[
\begin{align*}
e & \cup \mathcal{A} \cdot r \\
= & \{ (23); \mathcal{A} \text{ is uni-disjunctive} \} \\
e & \cup (\cup k : \mathcal{O} \cdot (\mathcal{A}^k \cdot e)) \\
= & \{ \text{definition of exponentiation; dummy } k := k - 1 \} \\
\mathcal{A}^0 \cdot e & \cup (\cup k : \mathcal{A}^k \cdot e) \\
= & \{ \text{domain split; } (23) \} \\
r.
\end{align*}
\]
Secondly, it must be shown that \( r \) satisfies (25). Using a property of \( \mathcal{U} \) and definition (23) of \( r \), we have, for arbitrary \( s \):

\[
 r \subseteq s \equiv (\forall k: 0 < k: \mathcal{A}^k \cdot e \subseteq s). \tag{B0}
\]

The right-hand side of (B0) can be proven with induction to \( k \), assuming the validity of \( e \cup \mathcal{A} \cdot s \subseteq s \). The base, \( \mathcal{A}^0 \cdot e = e \subseteq s \), follows trivially. For the induction step it follows that

\[
\begin{align*}
\mathcal{A}^{k+1} \cdot e &= \langle \text{definition of exponentiation} \rangle \\
&= \mathcal{A} \cdot (\mathcal{A}^k \cdot e) \\
&\subseteq \langle \mathcal{A} \text{ is monotonic; induction hypothesis } \mathcal{A}^k \cdot e \subseteq s \rangle \\
&\mathcal{A} \cdot s \\
&\subseteq \langle \text{assumption; property of } \mathcal{U} \rangle \\
&\quad s.
\end{align*}
\]

Finally, the uniqueness of the solution of (24) and (25) is to be shown. Let \( r_0 \) and \( r_1 \) both satisfy (24) and (25). Application of (25), for \( r_0 \), taking \( s \) to be \( r_1 \), we find \( r_0 \subseteq r_1 \), since \( r_0 \) satisfies (24), and hence \( e \cup \mathcal{A} \cdot r_0 \subseteq r_0 \). The reverse \( r_1 \subseteq r_0 \) follows analogously. Antisymmetry then implies \( r_0 = r_1 \).
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