# Loops under Strategies 

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#### Abstract

Most techniques to automatically disprove termination of term rewrite systems search for a loop. Whereas a loop implies nontermination for full rewriting, this is not necessarily the case if one considers rewriting under strategies. Therefore, in this paper we first generalize the notion of a loop to a loop under a given strategy. In a second step we present two novel decision procedures to check whether a given loop is a context-sensitive or an outermost loop. We implemented and successfully evaluated our method in the termination prover $\mathrm{T}_{\boldsymbol{T}} \mathrm{T}_{2}$.


## 1 Introduction

Termination is an important property of term rewrite systems (TRSs). Therefore, much effort has been spent on developing and automating powerful techniques for showing termination of TRSs. An important application area for these techniques is termination analysis of functional programs. Since the evaluation mechanism of functional languages is mainly term rewriting, one can transform functional programs into TRSs and prove termination of the resulting TRSs to conclude termination of the functional programs [6]. Although "full" rewriting does not impose any evaluation strategy, this approach is sound even if the underlying programming language has an evaluation strategy.

But in order to detect bugs in programs, it is at least as important to prove non-termination of programs or of the corresponding TRSs. Here, the evaluation strategy cannot be ignored, because a non-terminating TRS may still be terminating when considering the strategy. Thus, in order to disprove termination of programming languages with strategies, it is important to develop automated techniques to disprove termination of TRSs under strategies.

Only a few techniques for showing non-termination of TRSs have been introduced so far $[4,7,9,10,12]$. These techniques can be used to detect loops-a specific form of derivation which implies non-termination-and are successfully implemented in many tools (e.g., AProVE[5], Jambox [2], Matchbox [16], NTI [12], TORPA [17], $\mathrm{T}^{\top} \mathrm{T}_{2}$ [8]).

If one wants to prove non-termination under strategies then up to now there are two different approaches. The first one is to directly analyze the loops whether they also imply non-termination under a given strategy $\mathcal{S}$. This approach was successfully applied for the innermost strategy in [15] where a decision procedure was given to determine whether a loop is an innermost loop.

[^0]The second approach is to use a complete transformation $\tau_{\mathcal{S}}$ for strategy $\mathcal{S}$ such that $\mathcal{R}$ is terminating under $\mathcal{S}$ iff $\tau_{\mathcal{S}}(\mathcal{R})$ is (innermost) terminating. Then one first applies the transformation and then searches for a loop in $\tau_{\mathcal{S}}(\mathcal{R})$ afterwards. Here, the methods of [3] and [13,14] are applicable which can be used to disprove context-sensitive and outermost termination.

Although the second approach of using the transformations [3,13,14] seems to be a good solution to disprove context-sensitive and outermost termination, there are two main drawbacks. The first problem is a practical one. Often the loops of $\mathcal{R}$ are transformed into much longer loops in $\tau_{\mathcal{S}}(\mathcal{R})$ and hence, the search space for loops may become critical. And even more severe is the problem, that some loops of $\mathcal{R}$ are not even translated to loops in $\tau_{\mathcal{S}}(\mathcal{R})$ and hence, one even looses power if the search problem for loops is ignored.

Thus, there is still need to extend the first approach-to ensure or even decide that a given loop is a loop under strategies-to other strategies besides innermost. To this end, in this paper we first generalize the notion of a loop and an innermost loop to a loop under some arbitrary strategy. Then we develop two new decision procedures for context-sensitive loops and outermost loops.

The paper is structured as follows. In Sect. 2 we recapitulate the required notions of rewriting and generalize the notion of a loop for rewriting strategies. Moreover, we present a decision procedure for the question whether a given loop is a context-sensitive loop. Then in Sect. 3 we show how to formulate the same question for the outermost strategy as a set of matching problems. How these matching problems can be transformed to a simpler kind of problems-identity problems - is the content of Sect. 4. Afterwards, in Sect. 5 we provide a decision procedure for solvability of identity problems. All of our techniques have been implemented in the Tyrolean Termination Tool $2\left(\mathrm{~T}_{\mathrm{T}} \mathrm{T}_{2}\right)$ and the empirical results are presented in Sect. 6, before we conclude in Sect. 7.

## 2 Loops

We only regard finite signatures and TRSs and refer to [1] for the basics of rewriting. We use $\ell, r, s, t, u, \ldots$ for terms, $f, g, \ldots$ for function symbols, $x, y, \ldots$ for variables, $\sigma, \mu$ for substitutions, $i, j, k, n, m, o$ for natural numbers, $p, q, \ldots$ for positions where $\varepsilon$ is the root position, and $C, D, \ldots$ for contexts. Here, contexts are terms which contain exactly one hole $\square$. For contexts, the term $C[t]$ is like $C$ where $\square$ is replaced by $t$, i.e., $\square[t]=t$ and $f\left(s_{1}, \ldots, C, \ldots, s_{n}\right)[t]=$ $f\left(s_{1}, \ldots, C[t], \ldots, s_{n}\right)$. We write $\left.t\right|_{p}$ for the subterm of $t$ at position $p$, i.e., $\left.t\right|_{\varepsilon}=t$ and $\left.f\left(s_{1}, \ldots, s_{n}\right)\right|_{i p}=\left.s_{i}\right|_{p}$. The set of variables is denoted by $\mathcal{V}$.

Throughout this paper we assume a fixed TRS $\mathcal{R}$ and we write $t \rightarrow_{p} s$ if one can reduce $t$ to $s$ at position $p$ with $\mathcal{R}$, i.e., $t=C[\ell \sigma]$ and $s=C[r \sigma]$ for some $\ell \rightarrow r \in \mathcal{R}$, substitution $\sigma$, and context $C$ with $\left.C\right|_{p}=\square$. Here, the term $\ell \sigma$ is called a redex at position $p$. The reduction is an outermost reduction, written $t{ }^{\circ}{ }_{p} s$, iff $t$ contains no redex at a position $q$ above $p$ (written $q<p$ ). If the position is irrelevant we just write $\rightarrow$ or $\xrightarrow{\circ}$. The TRS $\mathcal{R}$ is non-terminating iff
there is an infinite derivation $t_{1} \rightarrow t_{2} \rightarrow \ldots$. It is outermost non-terminating iff there is such an infinite derivation using $\xrightarrow{\circ}$ instead of $\rightarrow$.

An obvious approach to disprove termination is to search for a loop, i.e., a derivation where the starting term $t$ is reduced to a term containing an instance of $t$, i.e., $t \rightarrow^{+} C[t \mu]$. The corresponding infinite derivation is

$$
t \rightarrow^{+} C[t \mu] \rightarrow^{+} C[C[t \mu] \mu] \rightarrow^{+} \ldots \rightarrow^{+} C[C[\ldots C[t \mu] \ldots \mu] \mu] \rightarrow^{+} \ldots
$$

where the derivation $t \rightarrow^{+} C[t \mu]$ is repeated over and over again. This infinite derivation $(\star)$ is obtained because $\rightarrow$ is closed under both substitutions and contexts, also known as stability and monotonicity.

However, in general it is not clear whether $(\star)$ also is an infinite derivation if one considers a specific evaluation strategy $\mathcal{S}$. If this is the case then and only then we speak of an $\mathcal{S}$-loop.

To formally define an $\mathcal{S}$-loop we first need to make the derivations within (*) precise. Therefore, we must represent terms like $C[C[\ldots C[t \mu] \ldots \mu] \mu]$ without using "...". Moreover, we must know the positions of the reductions since several strategies - like innermost, outermost, or context-sensitive only allow reductions at certain positions. To this end, we define the notion of a contextsubstitution which combines insertion into a context with the application of a substitution.

Definition 1 (Context-substitutions). A context-substitution is a pair $(C, \mu)$ consisting of a context $C$ and a substitution $\mu$. The $n$-fold application of $(C, \mu)$ to a term $t$, written $t(C, \mu)^{n}$ is defined as follows.

- $t(C, \mu)^{0}=t$
- $t(C, \mu)^{n+1}=C\left[t(C, \mu)^{n} \mu\right]$

From the definition it is obvious that in $t(C, \mu)^{n}$ the context $C$ is added $n$-times above $t$ and $t$ is instantiated by $\mu^{n}$. Note that also the added contexts are instantiated by $\mu$. For the term $t(C, \mu)^{3}$ this is illustrated in Fig. 1.

The following lemma shows that context-substitutions have similar properties to both contexts and substitutions.


Fig. 1. The term $t(C, \mu)^{3}$

## Lemma 1 (Properties of context-substitutions).

(i) $t(C, \mu)^{n} \mu=t \mu(C \mu, \mu)^{n}$.
(ii) $t(C, \mu)^{m}(C, \mu)^{n}=t(C, \mu)^{m+n}$.
(iii) If $\left.C\right|_{p}=\square$ then $\left.t(C, \mu)^{n}\right|_{p^{n}}=t \mu^{n}$.
(iv) Whenever $t \rightarrow_{q} s$ and $\left.C\right|_{p}=\square$ then $t(C, \mu)^{n} \rightarrow_{p^{n} q} s(C, \mu)^{n}$.

Proof. (i) We perform induction on $n$. If $n=0$ then

$$
t(C, \mu)^{n} \mu=t(C, \mu)^{0} \mu=t \mu=t \mu(C \mu, \mu)^{0}=t \mu(C \mu, \mu)^{n}
$$

Otherwise, $n=k+1$ and we obtain

$$
\begin{aligned}
t(C, \mu)^{n} \mu & =t(C, \mu)^{k+1} \mu \\
& =C\left[t(C, \mu)^{k} \mu\right] \mu \\
\text { (by ind.) } & =C\left[t \mu(C \mu, \mu)^{k}\right] \mu \\
& =C \mu\left[t \mu(C \mu, \mu)^{k} \mu\right] \\
& =t \mu(C \mu, \mu)^{k+1} \\
& =t \mu(C \mu, \mu)^{n}
\end{aligned}
$$

(ii) We perform induction on $n$. If $n=0$ then

$$
t(C, \mu)^{m}(C, \mu)^{n}=t(C, \mu)^{m}(C, \mu)^{0}=t(C, \mu)^{m}=t(C, \mu)^{m+n}
$$

Otherwise, $n=k+1$ and we obtain

$$
\begin{aligned}
t(C, \mu)^{m}(C, \mu)^{n} & =t(C, \mu)^{m}(C, \mu)^{k+1} \\
& =C\left[t(C, \mu)^{m}(C, \mu)^{k} \mu\right] \\
(\text { by ind. }) & =C\left[t(C, \mu)^{m+k} \mu\right] \\
& =t(C, \mu)^{m+k+1} \\
& =t(C, \mu)^{m+n}
\end{aligned}
$$

(iii) We perform induction on $n$. If $n=0$ then $\left.t(C, \mu)^{n}\right|_{p^{n}}=\left.t(C, \mu)^{0}\right|_{\varepsilon}=t=$ $t \mu^{0}=t \mu^{n}$. Otherwise, $n=k+1$ and we obtain

$$
\begin{aligned}
\left.t(C, \mu)^{n}\right|_{p^{n}} & =\left.t(C, \mu)^{k+1}\right|_{p p^{k}} \\
& =\left.C\left[t(C, \mu)^{k} \mu\right]\right|_{p p^{k}} \\
& =\left.t(C, \mu)^{k} \mu\right|_{p^{k}} \\
& =\left.t(C, \mu)^{k}\right|_{p^{k}} \mu \\
\text { (by ind.) } & =\left(t \mu^{k}\right) \mu \\
& =t \mu^{n} .
\end{aligned}
$$

(iv) We perform induction on $n$. If $n=0$ then

$$
t(C, \mu)^{n}=t \rightarrow_{p^{0} q} s=s(C, \mu)^{n}
$$

Otherwise, $n=k+1$ and by induction hypothesis we have $t(C, \mu)^{k} \rightarrow p^{k} q$ $s(C, \mu)^{k}$. Thus,

$$
t(C, \mu)^{n}=C\left[t(C, \mu)^{k} \mu\right] \rightarrow_{p p^{k} q} C\left[s(C, \mu)^{k} \mu\right]=s(C, \mu)^{k+1}=s(C, \mu)^{n}
$$

Here, property (i) is similar to the fact that $C[t] \mu=C \mu[t \mu]$, and (ii) expresses that context-substitutions can be combined just as substitutions where $\sigma^{m} \sigma^{n}=$ $\sigma^{m+n}$. Moreover, the property of contexts that $C[t]_{p}=t$ if $\left.C\right|_{p}=\square$ is extended in (iii), and finally stability and monotonicity of rewriting are used to show in (iv) that rewriting is closed under context-substitutions.

With the help of context-substitutions we can now describe the infinite derivation in $(\star)$ more concisely. Since $t \rightarrow^{+} C[t \mu]=t(C, \mu)$ we obtain

$$
t(C, \mu)^{0} \rightarrow^{+} t(C, \mu)(C, \mu)^{0}=t(C, \mu)^{1} \rightarrow^{+} \ldots \rightarrow^{+} t(C, \mu)^{n} \rightarrow^{+} \ldots
$$

Hence, the terms that occur during the derivation are precisely defined and for every $n$ the positions of the reductions are prefixed by an additional $p^{n}$ where $p$ is the position of the hole in $C$, cf. Lemma 1 (iv). In other words, every reduction takes place at the same position of the subterm $t \mu^{n}$ of $t(C, \mu)^{n}$.

Now it is natural to define that a derivation $t \rightarrow^{+} t(C, \mu)$ is called an $\mathcal{S}$-loop iff all steps in ( $\star \star$ ) respect the strategy $\mathcal{S} .{ }^{1}$

Definition 2 (S-loops). Let $\mathcal{S}$ be a strategy. A loop $t_{1} \rightarrow_{q_{1}} t_{2} \rightarrow q_{2} \ldots t_{n} \rightarrow q_{n}$ $t_{n+1}=t_{1}(C, \mu)$ with $\left.C\right|_{p}=\square$ is an $\mathcal{S}$-loop iff all reductions $t_{i}(C, \mu)^{m} \rightarrow_{p^{m} q_{i}}$ $t_{i+1}(C, \mu)^{m}$ respect the strategy $\mathcal{S}$ for all $1 \leq i \leq n$ and $m \geq 0$.

As a direct consequence of Def. 2 one can conclude that every $\mathcal{S}$-loop of a rewrite system $\mathcal{R}$ proves non-termination of $\mathcal{R}$ under strategy $\mathcal{S}$.

Example 1. We consider the $\operatorname{TRS} \mathcal{R}_{n}$ for arithmetic with $n$-bit numbers.

$$
\begin{align*}
\mathrm{p}(0) & \rightarrow 0 & (1) & \mathrm{plus}(0, y)
\end{align*} \rightarrow y
$$

Here, the last rule is used to model that an overflow occurred due to the $n$-bit restriction. We focus on the loops

$$
\begin{aligned}
t_{1}=\operatorname{minus}(x, \text { inf }) & \rightarrow \operatorname{minus}(x, \mathrm{~s}(\text { inf })) \rightarrow \mathrm{p}(\operatorname{minus}(x, \text { inf }))=C_{1}\left(t_{1}, \mu_{1}\right) \quad \text { and } \\
t_{2}=\operatorname{plus}(\text { inf }, y) & \rightarrow \operatorname{plus}(\mathrm{s}(\text { inf }), y) \rightarrow \mathrm{s}(\operatorname{plus}(\text { inf }, y))=C_{2}\left(t_{2}, \mu_{2}\right)
\end{aligned}
$$

where $\mu_{1}=\mu_{2}=\{ \}, C_{1}=\mathrm{p}(\square)$, and $C_{2}=\mathrm{s}(\square)$. Here, the first loop is an outermost loop, but the second one is not. The reason for the latter is that in every iteration one more s is created. Hence, this will lead to a redex w.r.t. Rule (9). Note that this example will be hard to handle with the transformational approaches: using [14] creates a TRS where all infinite reductions are non-looping and [13] is not even applicable due to non-left-linearity of Rule (4).

[^1]Note that a loop is not only determined by the precise derivation $t_{1} \rightarrow^{+} t_{n+1}$, but also by the specific context $C$ that is used. To see this consider the TRS $\mathcal{R}=\{\mathrm{a} \rightarrow \mathrm{f}(\mathrm{a}, \mathrm{a}), \mathrm{f}(\mathrm{f}(x, y), z) \rightarrow \mathrm{b}\}$. Then $\mathrm{a} \rightarrow \mathrm{f}(\mathrm{a}, \mathrm{a})$ is a looping derivation. For $C=\mathrm{f}(\mathrm{a}, \square)$ this is an outermost loop. However, for $C^{\prime}=\mathrm{f}(\square, \mathrm{a})$ we do not obtain an outermost loop since $f(a, a) \xrightarrow{\circ} \mathrm{f}(\mathrm{f}(\mathrm{a}, \mathrm{a}), \mathrm{a}) \xrightarrow{q} \mathrm{f}(\mathrm{f}(\mathrm{f}(\mathrm{a}, \mathrm{a}), \mathrm{a}), \mathrm{a})$. Hence, the choice of $C$ is essential.

For automatic non-termination analysis under strategies one main question is whether a given loop is an $\mathcal{S}$-loop, i.e., whether the loop implies non-termination even under strategy $\mathcal{S}$. In [15] it was already shown that this question is decidable for innermost loops. ${ }^{2}$ There, one has the problem that innermost rewriting is not stable, although it is monotonic.

In context-sensitive rewriting [11] we have the inverse situation: first, stability is given whereas monotonicity is absent. And second, whereas the decision procedure for innermost loops is quite involved and already known, for contextsensitive loops we can present a novel decision procedure that is rather straightforward, but nevertheless important.

Theorem 1 (Deciding context-sensitive-loops). A loop $t \rightarrow^{+} C[t \mu]$ is a context-sensitive loop (using replacement map $\nu$ ) iff both the derivation $t \rightarrow^{+}$ $C[t \mu]$ respects the context-sensitive strategy and the hole in $C$ is at a $\nu$-replacing position.

Proof. Let $t_{1} \rightarrow_{q_{1}} t_{2} \rightarrow q_{2} \ldots t_{n} \rightarrow q_{n} t_{n+1}=t_{1}(C, \mu)$ be a loop where $\left.C\right|_{p}=\square$. Then the following statements are all equivalent.

- the loop is a context-sensitive loop
- all $t_{i}(C, \mu)^{m} \rightarrow_{p^{m} q_{i}} t_{i+1}(C, \mu)^{m}$ are context-sensitive reductions
- all $p^{m} q_{i}$ are $\nu$-replacing positions of $t_{i}(C, \mu)^{m}$
- $p$ is a $\nu$-replacing position of $C$ and each $q_{i}$ is a $\nu$-replacing position of $t_{i} \mu^{m}$
- the hole in $C$ is at a $\nu$-replacing position and the derivation $t_{1} \rightarrow^{+} t_{1}(C, \mu)$ is a context-sensitive derivation

In the rest of this paper we consider the outermost strategy. As main result we develop a decision procedure for the question whether a given loop is an outermost loop. Note that for outermost rewriting neither stability nor monotonicity are given. To see this consider the TRS $\mathcal{R}=\{\mathrm{a} \rightarrow \mathrm{a}, \mathrm{f}(x) \rightarrow x, \mathrm{~g}(\mathrm{f}(\mathrm{a})) \rightarrow \mathrm{a}\}$. Then $\mathrm{a} \xrightarrow{\circ} \mathrm{a}$, but $\mathrm{f}(\mathrm{a}) \underset{\rightarrow}{\nrightarrow} \mathrm{f}(\mathrm{a})$. Moreover, $\mathrm{g}(\mathrm{f}(x)) \xrightarrow{\circ} \mathrm{g}(x)$, but $\mathrm{g}(\mathrm{f}(\mathrm{a})) \underset{\rightarrow}{q} \mathrm{~g}(\mathrm{a})$.

The problem of missing stability was already present for innermost loops. Therefore, many techniques of [15] for innermost loops can be reused for outermost loops, too. However, to handle the missing monotonicity of outermost rewriting we have to extend these techniques by an additional context. And these contexts will require significant extensions of the techniques of [15] and are not so easy to treat as in the context-sensitive case.

[^2]
## 3 Deciding Outermost Loops

Recall the definition of an outermost reduction. An outermost reduction of $t$ at position $p$ requires that there is no redex at a position $q$ above $p$, i.e., all subterms $\left.t\right|_{q}$ with $q<p$ must not be matched by some left-hand side of a rule in $\mathcal{R}$. Hence, the question of an outermost reduction can be formulated as a question of matching.

However, we do not have to consider a single outermost reduction but we want to know whether each reduction of a term $t(C, \mu)^{m}$ at position $p^{m} q$ is an outermost reduction (where $\left.C\right|_{p}=\square$ and $q \in \mathcal{P o s}(t)$ ). Looking at Fig. 1 one sees that there are two different cases how to obtain a subterm at a position above $p^{m} q$ that is matched by some left-hand side. First, the subterm may be a subterm of $t \mu^{m}$. Or otherwise, the subterm starts within the context. For the former case we can reuse the so called matching problems [15, Def. 12] and for the latter we need an extended version of matching problems containing contexts.

Definition 1 ((Extended) matching problems). A matching problem is a pair $(\mathcal{M}, \mu)$ where $\mathcal{M}$ is a set of pairs of terms $s \gtrdot \ell$. It is solvable iff there is a solution $(k, \sigma)$ such that for all $s \gtrdot \ell \in \mathcal{M}$ the equation $s \mu^{k}=\ell \sigma$ is satisfied.

An extended matching problem is a quintuple $(D \gtrdot \ell, C, t, \mathcal{M}, \mu)$. It is solvable iff there is a solution $(n, k, \sigma)$ such that the equation $D\left[t(C, \mu)^{n}\right] \mu^{k}=\ell \sigma$ is satisfied and $(k, \sigma)$ is a solution to the matching problem $(\mathcal{M}, \mu)$.

To simplify presentation we write $(D \gtrdot \ell, C, t, \mu)$ instead of $(D \gtrdot \ell, C, t, \varnothing, \mu)$ and we write $(s \gtrdot \ell, \mu)$ instead of $(\{s \gtrdot \ell\}, \mu)$. Moreover, we use the notion "matching problem" also for extended matching problems.

To check whether $t(C, \mu)^{m}$ has a redex above position $p^{m} q$ one can now construct a set of initial matching problems. Essentially, one considers matching problems for the subterms of $t$ above $q$. Additionally, for each subterm of $t(C, \mu)^{m}$ that starts with a subcontext $\left.C\right|_{p^{\prime}}$ of $C$, we build an extended matching problem.

Definition 2 (Initial matching problems). Let $t \rightarrow_{q} u$ be a reduction and $(C, \mu)$ be a context-substitution with $\left.C\right|_{p}=\square$. Then the following initial matching problems are created for this reduction and context-substitution.

- $\left(\left.t\right|_{p^{\prime}} \gtrdot \ell, \mu\right)$ for each $\ell \rightarrow r \in \mathcal{R}$ and $p^{\prime}<q$
- $\left(\left.C\right|_{p^{\prime}} \gtrdot \ell, C \mu, t \mu, \mu\right)$ for each $\ell \rightarrow r \in \mathcal{R}$ and $p^{\prime}<p$

Example 1. Consider the loop $t_{1}=\operatorname{minus}(x, \inf ) \rightarrow_{2} \operatorname{minus}(x, \mathrm{~s}(\mathrm{inf}))=t_{2} \rightarrow_{\varepsilon}$ $\mathrm{p}(\operatorname{minus}(x, \inf ))=t_{1}(C, \mu)$ of Ex. 1 where $C=\mathrm{p}(\square)$ and $\mu=\{ \}$. For the second reduction at root position we only build the extended matching problems $M P_{1 \ell}=(\mathrm{p}(\square) \gtrdot \ell, \mathrm{p}(\square)$, minus $(x, \mathrm{~s}(\mathrm{inf})), \mu)$ for all left-hand sides $\ell$ of $\mathcal{R}$. For the first reduction we obtain the similar extended matching problems $M P_{2 \ell}=$ $(\mathrm{p}(\square) \gtrdot \ell, \mathrm{p}(\square)$, minus $(x, \inf ), \mu)$, but additionally we also get the matching problems $M P_{3 \ell}=(\operatorname{minus}(x$, inf $) \gtrdot \ell, \mu)$.

The following theorem states that we have setup the right initial matching problems. If we consider all initial problems of all reductions $t_{i} \rightarrow q_{i} t_{i+1}$ of a
loop, then solvability of one of these problems is equivalent to the property that the loop is not outermost.
Theorem 1 (Outermost loops and matching problems). Let $t \rightarrow_{q} u$ and $(C, \mu)$ be given such that $\left.C\right|_{p}=\square$. All reductions $t(C, \mu)^{m} \rightarrow_{p^{m} q} u(C, \mu)^{m}$ are outermost iff none of the initial matching problems for $t \rightarrow{ }_{q} u$ and $(C, \mu)$ is solvable.

Proof. We first prove that any solvable initial matching problem shows that there is at least one reduction of $t(C, \mu)^{m}$ at position $p^{m} q$ which is not an outermost reduction. There are two cases. First, if $\left(\left.t\right|_{p^{\prime}} \gtrdot \ell, \mu\right)$ is solvable, then there is a solution $(k, \sigma)$ such that $\left.t\right|_{p^{\prime}} \mu^{k}=\ell \sigma$. Then $\left.t(C, \mu)^{k}\right|_{p^{k} p^{\prime}}=\left.t \mu^{k}\right|_{p^{\prime}}=\left.t\right|_{p^{\prime}} \mu^{k}=\ell \sigma$ shows that there is a redex in $t(C, \mu)^{k}$ above $p^{k} q$. Thus, $t(C, \mu)^{k} \rightarrow_{p^{k} q} u(C, \mu)^{k}$ is not an outermost reduction.

Otherwise, $\left(\left.C\right|_{p^{\prime}} \gtrdot \ell, C \mu, t \mu, \mu\right)$ is solvable. Hence, there is a solution $(n, k, \sigma)$ such that $\left.C\right|_{p^{\prime}}\left[t \mu(C \mu, \mu)^{n}\right] \mu^{k}=\ell \sigma$. Here, we show that the term $t(C, \mu)^{n+1+k}$ has a redex at position $p^{k} p^{\prime}$ which is above position $p^{n+1+k} q$ :

$$
\begin{array}{rlrl}
\left.t(C, \mu)^{n+1+k}\right|_{p^{k} p^{\prime}} & =\left.t(C, \mu)^{n}(C, \mu)(C, \mu)^{k}\right|_{p^{k} p^{\prime}} & =\left.t(C, \mu)^{n}(C, \mu) \mu^{k}\right|_{p^{\prime}} \\
& =\left.C\left[t(C, \mu)^{n} \mu\right] \mu^{k}\right|_{p^{\prime}} & & =\left.C\right|_{p^{\prime}}\left[t(C, \mu)^{n} \mu\right] \mu^{k} \\
& =\left.C\right|_{p^{\prime}}\left[t \mu(C \mu, \mu)^{n}\right] \mu^{k} & & =\ell \sigma
\end{array}
$$

For the other direction we show that if some reduction $t(C, \mu)^{m}$ at position $p^{m} q$ is not an outermost reduction, then one of the initial matching problems must be solvable. So suppose, the reduction of $t(C, \mu)^{m}$ is not outermost. Then there must be some position $q^{\prime}<p^{m} q$ such that the corresponding subterm $\left.t(C, \mu)^{m}\right|_{q^{\prime}}$ is a redex $\ell \sigma$. Again, there are two cases.

First, if $q^{\prime} \geq p^{m}$ then $q^{\prime}=p^{m} p^{\prime}$ where $p^{\prime}<q$ as $q^{\prime}<p^{m} q$. Hence, $\ell \sigma=$ $\left.t(C, \mu)^{m}\right|_{q^{\prime}}=\left.t(C, \mu)^{m}\right|_{p^{m} p^{\prime}}=\left.t \mu^{m}\right|_{p^{\prime}}=\left.t\right|_{p^{\prime}} \mu^{m}$. Thus, the initial matching problem $\left(\left.t\right|_{p^{\prime}} \gtrdot \ell, \mu\right)$ has the solution $(m, \sigma)$.

In the other case $q^{\prime}<p^{m}$. Thus, we can split the position $q^{\prime}$ into $p^{k} p^{\prime}$ where $k<m$ and $p^{\prime}<p$. Moreover, there must be some $n \in \mathbb{N}$ that satisfies $m=$ $n+1+k$. We conclude

$$
\begin{aligned}
\ell \sigma & =\left.t(C, \mu)^{m}\right|_{q^{\prime}} & & =\left.t(C, \mu)^{n+1+k}\right|_{p^{k} p^{\prime}} \\
& =\left.t(C, \mu)^{n}(C, \mu)(C, \mu)^{k}\right|_{p^{k} p^{\prime}} & =\left.t(C, \mu)^{n}(C, \mu) \mu^{k}\right|_{p^{\prime}} & =\left.C\left[t \mu(C \mu, \mu)^{n}\right] \mu^{k}\right|_{p^{\prime}} \\
& =\left.C\left[t(C, \mu)^{n} \mu\right] \mu^{k}\right|_{p^{\prime}} & & =\left.C\right|_{p^{\prime}}\left[t \mu(C \mu, \mu)^{n}\right] \mu^{k} .
\end{aligned}
$$

Thus, the initial matching problem $\left(\left.C\right|_{p^{\prime}} \gtrdot \ell, C \mu, t \mu, \mu\right)$ is solvable.
Note that whenever $C=\square$ then there is no initial matching problem which is an extended matching problem. Hence, by Thm. 1 one can already decide whether a loop $t \rightarrow^{+} t \mu$ is an outermost loop by using the techniques of [15] to decide solvability of matching problems. For example it can be detected that all matching problems $M P_{3 \ell}$ of Ex. 1 are not solvable.

However, in the general case we also generate extended matching problems. Therefore, in the next section we develop a novel decision procedure for solvability of extended matching problems like $M P_{1 \ell}$ and $M P_{2 \ell}$ of Ex. 1 .

## 4 Deciding Solvability of Extended Matching Problems

Since extended matching problems are only generated if $C \neq \square$, in the following sections we always assume that $C \neq \square$. We take a similar approach to [15] where we transform each matching problem into $\top$, $\perp$, or into solved form. Here $\top$ and $\perp$ represent solvability and non-solvability. And if a matching problem is in solved form then often solvability can immediately be decided. We explain all transformation rules in detail directly after the following definition.

Definition 1 (Transformation of extended matching problems). Let $M P=$ $\left(D \gtrdot \ell_{0}, C, t, \mathcal{M}, \mu\right)$ be an extended matching problem where $\mathcal{M}=\left\{s_{1} \gtrdot \ell_{1}, \ldots, s_{m} \gtrdot\right.$ $\left.\ell_{m}\right\}$. Then MP is in solved form iff each $\ell_{i}$ is a variable. Let $\mathcal{V}_{i n c r}=\{x \in \mathcal{V} \mid$ $\left.\exists n: x \mu^{n} \notin \mathcal{V}\right\}$ be the set of increasing variables.

We define a relation $\Rightarrow$ which simplifies extended matching problems that are not in solved form. So, let $\ell_{j}=f\left(\ell_{1}^{\prime}, \ldots, \ell_{m^{\prime}}^{\prime}\right)$.
(i) $M P \Rightarrow\left(D_{i^{\prime}} \gtrdot \ell_{i^{\prime}}^{\prime}, C, t, \mathcal{M} \cup\left\{t_{i} \gtrdot \ell_{i}^{\prime} \mid 1 \leq i \leq m^{\prime}, i \neq i^{\prime}\right\}\right.$, $\left.\mu\right)$ if $j=0$ and $D=$ $f\left(t_{1}, \ldots, D_{i^{\prime}}, \ldots, t_{m^{\prime}}\right)$.
(ii) $M P \Rightarrow\left(D \gtrdot \ell_{0}, C, t,\left(\mathcal{M} \backslash\left\{s_{j} \gtrdot \ell_{j}\right\}\right) \cup\left\{t_{i} \gtrdot \ell_{i}^{\prime} \mid 1 \leq i \leq m^{\prime}\right\}, \mu\right)$ if $j>0$ and $s_{j}=f\left(t_{1}, \ldots, t_{m^{\prime}}\right)$.
(iii) $M P \Rightarrow \perp$ if $j=0$ and $D=g(\ldots)$ where $f \neq g$.
(iv) $M P \Rightarrow \perp$ if $j>0$ and $s_{j}=g(\ldots)$ where $f \neq g$.
(v) $M P \Rightarrow \perp$ if $j>0$ and $s_{j} \in \mathcal{V} \backslash \mathcal{V}_{\text {incr }}$.
(vi) $M P \Rightarrow\left(D \mu \gtrdot \ell_{0}, C \mu, t \mu,\left\{s_{i} \mu \gtrdot \ell_{i} \mid 1 \leq i \leq m\right\}, \mu\right)$ if $j>0$ and $s_{j} \in \mathcal{V}_{\text {incr }}$.
(vii) $M P \Rightarrow \top$ if $j=0, D=\square$, and $\left(\mathcal{M} \cup\left\{t \gtrdot \ell_{0}\right\}, \mu\right)$ is solvable.
(viii) $M P \Rightarrow\left(C \gtrdot \ell_{0}, C \mu, t \mu, \mathcal{M}, \mu\right)$ if $j=0, D=\square$, and $\left(\mathcal{M} \cup\left\{t \gtrdot \ell_{0}\right\}, \mu\right)$ is not solvable.

Recall that $M P=\left(D \gtrdot \ell_{0}, C, t,\left\{s_{1} \gtrdot \ell_{1}, \ldots, s_{m} \gtrdot \ell_{m}\right\}, \mu\right)$ is solvable iff there is a solution $(n, k, \sigma)$ such that $D\left[t(C, \mu)^{n}\right] \mu^{k}=\ell_{0} \sigma$ and $s_{i} \mu^{k}=\ell_{i} \sigma$ for all $1 \leq$ $i \leq m$. Hence, whenever $D \neq \square$ or $s_{i} \notin \mathcal{V}$ then one can perform a decomposition (Rules (i) and (ii)) or detect a clash (Rules (iii) and (iv)) as in a standard matching algorithm.

If $s_{j}=x$ is a non-increasing variable then $s_{j} \mu^{k}$ will always be a variable. Thus, Rule (v) correctly returns $\perp$. But if $s_{j}=x$ is an increasing variable then there might be a solution if $k>0$. Hence, one can just apply $\mu$ once on the whole matching problem using Rule (vi). Note that in the result of Rule (vi) both $t$ and $C$ are also instantiated. This reflects the property of context-substitutions that $t(C, \mu)^{n} \mu=t \mu(C \mu, \mu)^{n}$, cf. Lemma 1.

Whereas Rule (v) and a simplified version of Rule (vi) have already been present in [15], here we also need two additional rules to handle contexts. Note that for $n=0$ and $D=\square$ the term $D\left[t(C, \mu)^{n}\right] \mu^{k}$ is just $t \mu^{k}$ and thus, one only has to consider a non-extended matching problem. Now, in Rules (vii) and (viii) there is a case distinction whether this non-extended matching problem is solvable, i.e., whether $n=0$ yields a solution or not. If it is solvable then also a solution of $M P$ is found and Rule (vii) correctly returns $T$. If it is not possible
then there is only one way to continue: apply the context-substitution at least once, and this is exactly what Rule (viii) does.

Before we formally state the soundness of the transformation rules in Thm. 1 we illustrate their application on the extended matching problems of Ex. 1.

Example 1. We first consider $M P_{1 \ell}=(\mathrm{p}(\square) \gtrdot \ell, \mathrm{p}(\square)$, minus $(x, \mathrm{~s}(\mathrm{inf})), \mu)$. If $\ell$ is not one of the left-hand sides $\mathrm{p}(0)$ or $\mathrm{p}(\mathrm{s}(x))$ then $\perp$ is obtained by Rule (iii). If one considers $\ell=\mathrm{p}(0)$ then $(\mathrm{p}(\square) \gtrdot \mathrm{p}(0), \mathrm{p}(\square)$, minus $(x, \mathrm{~s}($ inf $)), \mu) \Rightarrow$ $(\square \gtrdot 0, \mathrm{p}(\square)$, minus $(x, \mathrm{~s}($ inf $)), \mu)$ by Rule (i). And as $(\operatorname{minus}(x, \mathrm{~s}(\mathrm{inf})) \gtrdot 0, \mu$ ) is not solvable, Rule (viii) yields $(\mathrm{p}(\square) \gtrdot 0, \mathrm{p}(\square)$, minus $(x, \mathrm{~s}(\mathrm{inf})), \mu)$. Finally, an application of Rule (iii) returns $\perp$ and thereby shows that the matching problem is not solvable. Since the transformation for $\ell=\mathrm{p}(\mathrm{s}(x))$ also results in $\perp$, we have detected that none of the matching problems $M P_{1 \ell}$ is solvable.

A similar transformation shows that none of the matching problems $M P_{2 \ell}$ is solvable. Hence, the loop of Ex. 1 is an outermost loop.

## Theorem 1 (Soundness and termination of the transformation rules).

(i) If $M P \Rightarrow \perp$ then $M P$ is not solvable.
(ii) If $M P \Rightarrow \top$ then $M P$ is solvable.
(iii) If $M P \Rightarrow M P^{\prime}$ then $M P$ is solvable iff $M P^{\prime}$ is solvable.
(iv) The relation $\Rightarrow$ is terminating and confluent. ${ }^{3}$

Proof. (i) There are three cases. If $j=0$ and $D=g(\ldots)$ then $M P$ is not solvable since $D\left[t(C, \mu)^{n}\right] \mu^{k}=g(\ldots) \mu^{k}=g(\ldots) \neq f(\ldots)=\ell_{0} \sigma$. Similarly, if $j>0$ and $s_{j}=g(\ldots)$ then $s_{j} \mu^{k}=g(\ldots) \neq f(\ldots)=\ell_{j} \sigma$ shows that $M P$ is not solvable. In the last case, $j>0$ and $s_{j}=x \in \mathcal{V} \backslash \mathcal{V}_{\text {incr }}$. Then, $s_{j} \mu^{k}=x \mu^{k} \neq f(\ldots)=\ell \sigma$ where the inequality is due to the fact that $x \mu^{k}$ will always be a variable. The reason is that $x \notin \mathcal{V}_{\text {incr }}$.
(ii) If $M P \Rightarrow \top$ then $j=0, D=\square$, and $\left(\mathcal{M} \cup\left\{t \gtrdot \ell_{0}\right\}, \mu\right)$ is solvable, i.e., there is some solution $(k, \sigma)$ such that for all $\left(u_{i}, \ell_{i}\right) \in \mathcal{M} \cup\left\{t \gtrdot \ell_{0}\right\}$ we have $u_{i} \mu^{k}=\ell_{i} \sigma$. Here, we identify $u_{0}=t$ and $u_{i}=s_{i}$ for $i>0$. We prove that $(n=0, k, \sigma)$ is a solution for $M P$. To this end, first observe that for $i>0$ we directly have $s_{i} \mu^{k}=u_{i} \mu_{k}=\ell_{i} \sigma$. Moreover, we have $D\left[t(C, \mu)^{n}\right] \mu^{k}=\square\left[u_{0}(C, \mu)^{0}\right] \mu^{k}=u_{0} \mu^{k}=\ell_{0} \sigma$.
(iii) We consider all rules separately and start with Rule (i), so let $j=0$ and $D=f\left(t_{1}, \ldots, D_{i^{\prime}}, \ldots, t_{m^{\prime}}\right)$. We show that $(n, k, \sigma)$ is a solution for $M P$ iff it is a solution for $M P^{\prime}=\left(D_{i^{\prime}} \gtrdot \ell_{i^{\prime}}^{\prime}, C, t, \mathcal{M} \cup\left\{t_{i} \gtrdot \ell_{i}^{\prime} \mid 1 \leq i \leq m^{\prime}, i \neq i^{\prime}\right\}, \mu\right)$. As $\mathcal{M}$ occurs in both $M P$ and $M P^{\prime}$ and as $\mu, k$, and $\sigma$ are unchanged we only have to consider the remaining equalities, i.e., the equality $D\left[t(C, \mu)^{n}\right] \mu^{k}=$ $\ell_{0} \sigma$ for $M P$ and the equalities $D_{i^{\prime}}\left[t(C, \mu)^{n}\right] \mu^{k}=\ell_{i}^{\prime} \sigma$ and $t_{i} \mu^{k}=\ell_{i}^{\prime} \sigma$ for each $i \neq i^{\prime}, 1 \leq i \leq m^{\prime}$. But since

$$
\begin{aligned}
D\left[t(C, \mu)^{n}\right] \mu^{k} & =f\left(t_{1}, \ldots, D_{i^{\prime}}, \ldots, t_{m^{\prime}}\right)\left[t(C, \mu)^{n}\right] \mu^{k} \\
& =f\left(t_{1}, \ldots, D_{i^{\prime}}\left[t(C, \mu)^{n}\right], \ldots, t_{m^{\prime}}\right) \mu^{k} \\
& =f\left(t_{1} \mu^{k}, \ldots, D_{i^{\prime}}\left[t(C, \mu)^{n}\right] \mu^{k}, \ldots, t_{m^{\prime}} \mu^{k}\right)
\end{aligned}
$$

[^3]and $\ell_{0} \sigma=\ell_{j} \sigma=f\left(\ell_{1}^{\prime}, \ldots, \ell_{m^{\prime}}^{\prime}\right) \sigma=f\left(\ell_{1}^{\prime} \sigma, \ldots, \ell_{m^{\prime}}^{\prime} \sigma\right)$ we conclude that $D\left[t(C, \mu)^{n}\right] \mu^{k}=\ell_{0} \sigma$ iff both $D_{i^{\prime}}\left[t(C, \mu)^{n}\right] \mu^{k}=\ell_{i}^{\prime} \sigma$ and $t_{i} \mu^{k}=\ell_{i}^{\prime} \sigma$ for each $i \neq i^{\prime}, 1 \leq i \leq m^{\prime}$.
The reasoning for Rule (ii) is similar to the one for Rule (i). Let $j>0$ and $s_{j}=f\left(t_{1}, \ldots, t_{m^{\prime}}\right)$. We show that $(n, k, \sigma)$ is a solution for MP iff it is a solution for $M P^{\prime}=\left(D \gtrdot \ell_{0}, C, t,\left(\mathcal{M} \backslash\left\{s_{j} \gtrdot \ell_{j}\right\}\right) \cup\left\{t_{i} \gtrdot \ell_{i}^{\prime} \mid 1 \leq i \leq m^{\prime}\right\}, \mu\right)$. Since $\mathcal{M} \backslash\left\{s_{j} \gtrdot \ell_{j}\right\}$ occurs in both $M P$ and $M P^{\prime}$ and since $D, \mu, n, k$, and $\sigma$ are unchanged we only have to consider the remaining equalities, i.e., the equality $s_{j} \mu^{k}=\ell_{j} \sigma$ for $M P$ and the equalities $t_{i} \mu^{k}=\ell_{i}^{\prime} \sigma$ for each $1 \leq i \leq m^{\prime}$. But since $s_{j} \mu^{k}=f\left(t_{1}, \ldots, t_{m^{\prime}}\right) \mu^{k}=f\left(t_{1} \mu^{k}, \ldots, t_{m^{\prime}} \mu^{k}\right)$ and $\ell_{j} \sigma=f\left(\ell_{1}^{\prime}, \ldots, \ell_{m^{\prime}}^{\prime}\right) \sigma=f\left(\ell_{1}^{\prime} \sigma, \ldots, \ell_{m^{\prime}}^{\prime} \sigma\right)$ we know that $s_{j} \mu^{k}=\ell_{j} \sigma$ iff $t_{i} \mu^{k}=$ $\ell_{i}^{\prime} \sigma$ for each $1 \leq i \leq m^{\prime}$.
Next, we consider Rule (vi). First observe, that whenever there is a solution $(n, k, \sigma)$ for $M P$ then there also is a solution with $k+1$, namely $(n, k+$ $1, \sigma \mu)$. The reason is that $D\left[t(C, \mu)^{n}\right] \mu^{k}=\ell_{0} \sigma$ implies $D\left[t(C, \mu)^{n}\right] \mu^{k+1}=$ $D\left[t(C, \mu)^{n}\right] \mu^{k} \mu=\ell_{0} \sigma \mu$ and similarly, $s_{i} \mu^{k}=\ell_{i} \sigma$ implies $s_{i} \mu^{k+1}=\ell_{i} \sigma \mu$.
Now, $\left(D \mu \gtrdot \ell_{0}, C \mu, t \mu,\left\{s_{i} \mu \gtrdot \ell_{i} \mid 1 \leq i \leq m\right\}, \mu\right)$ is solvable iff there is a solution $(n, k, \sigma)$ such that $D \mu\left[t \mu(C \mu, \mu)^{n}\right] \mu^{k}=\ell_{0} \sigma$ and $s_{i} \mu^{k+1}=\ell_{i} \sigma$ for each $1 \leq i \leq m$. But since
\[

$$
\begin{aligned}
D \mu\left[t \mu(C \mu, \mu)^{n}\right] \mu^{k} & =D \mu\left[t(C, \mu)^{n} \mu\right] \mu^{k} \\
& =D\left[t(C, \mu)^{n}\right] \mu \mu^{k} \\
& =D\left[t(C, \mu)^{n}\right] \mu^{k+1}
\end{aligned}
$$
\]

this is equivalent to solvability of $M P$ by the above observation.
Finally we consider Rule (viii), where $M P^{\prime}=\left(\mathcal{M} \cup\left\{t \gtrdot \ell_{0}\right\}, \mu\right)$ is not solvable, $j=0$, and $D=\square$. First observe, that since $M P^{\prime}$ is not solvable, there cannot be a solution of $M P$ where $n=0$. The reason is that if $(n=0, k, \sigma)$ would be a solution then $s_{j} \mu^{k}=\ell^{\prime} \sigma$ for each $1 \leq j \leq m$. Moreover, $t \mu^{k}=\square\left[t(C, \mu)^{0}\right] \mu^{k}=D\left[t(C, \mu)^{n}\right] \mu^{k}=\ell_{0} \sigma$. Thus, $(k, \sigma)$ would be a solution to the matching problem $M P^{\prime}$ in contradiction to the requirement that $M P^{\prime}$ is not solvable.
Using this observation we prove that $M P^{\prime \prime}=\left(C \gtrdot \ell_{0}, C \mu, t \mu, \mathcal{M}, \mu\right)$ is solvable iff $M P$ is solvable. By the above observation, $M P$ is solvable iff there is a solution $(n+1, k, \sigma)$ and we show that this is equivalent to the existence of a solution $(n, k, \sigma)$ for $M P^{\prime \prime}$. The reason is that since $k, \sigma, \mu$, and $\mathcal{M}$ are unchanged, we know that $k, \sigma$ is a solution for $(\mathcal{M}, \mu)$ as part of $M P$ iff it is a solution for $(\mathcal{M}, \mu)$ as part of $M P^{\prime \prime}$. And moreover, $D\left[t(C, \mu)^{n+1}\right] \mu^{k}=$ $\square\left[t(C, \mu)^{n+1}\right] \mu^{k}=t(C, \mu)^{n+1} \mu^{k}=C\left[t(C, \mu)^{n} \mu\right] \mu^{k}=C\left[t \mu(C \mu, \mu)^{n}\right] \mu^{k}$. So, the equalities concerning the contexts $D=\square$ of $M P$ and $C$ of $M P^{\prime \prime}$ are identical.
(iv) To show termination of $\Rightarrow$ first note that an application of one of the Rules (iii), (iv), (v), or (vii) obviously yields termination. For the remaining rules note that no transformation rule increases the size of the terms $\ell_{i}$. Thus, Rules (i) and (ii) can only be applied finitely often. But since every sequence of transformations with Rule (vi) eventually triggers an application of Rule
(ii) or (iv), Rule (vi) cannot be used infinitely often either. Similarly, each application of Rule (viii) triggers an application of Rule (i) or (iii) and thus, Rule (viii) can only be used finitely many times.

Confluence is due to the following two reasons. First, once $j$ has been chosen there is only one rule to apply since the rules are non-overlapping.

To show that the choice of $j$ is uncritical we perform a large but simple case distinction. Note that it cannot happen that $M P$ is reduced to both $\perp$ and $\top$ due to the previous results (i) and (ii) of this theorem. Next, we consider that $M P \Rightarrow \perp$ by Rule (iii) and $M P \Rightarrow M P^{\prime}$ using some other index $j^{\prime} \neq 0$. But then regardless of the chosen $j^{\prime}$ one can afterwards apply Rule (iii) again to reduce $M P^{\prime} \Rightarrow \perp$. Here one needs the fact, that if $D=g(\ldots)$ has a conflict with $\ell_{0}=f(\ldots)$ then $D \mu=g(\ldots)$ still has this conflict. Similarly, if $M P \Rightarrow \perp$ by Rule (iv) and $M P \Rightarrow M P^{\prime}$ then again $M P^{\prime} \Rightarrow \perp$ by Rule (iv). And also if $M P \Rightarrow \perp$ by Rule (v) then for all $M P \Rightarrow M P^{\prime}$ we know that $M P^{\prime} \Rightarrow \perp$ by Rule (v). Here, one needs the observation that $s_{j} \in \mathcal{V} \backslash \mathcal{V}_{\text {incr }}$ implies $s_{j} \mu \in \mathcal{V} \backslash \mathcal{V}_{\text {incr }}$. Next, we consider the case that $M P \Rightarrow \top$ by Rule (vii) and $M P \Rightarrow M P^{\prime}$. If one used Rule (vi) to obtain $M P^{\prime}$ then $M P^{\prime}=\left(\square \gtrdot \ell_{0}, C \mu, t \mu,\left\{s_{i} \mu \gtrdot \ell_{i} \mid 1 \leq i \leq m\right\}, \mu\right)$ and one can apply Rule (vii) afterwards to reduce $M P^{\prime} \Rightarrow \top$. This is possible since the resulting matching problem $\left(\left\{t \mu \gtrdot \ell_{0}\right\} \cup\left\{s_{i} \mu \gtrdot \ell_{i} \mid 1 \leq i \leq m\right\}, \mu\right)$ is satisfiable iff $\left(\left\{t \gtrdot \ell_{0}\right\} \cup\left\{s_{i} \gtrdot \ell_{i} \mid 1 \leq i \leq m\right\}, \mu\right)$ is satisfiable. Otherwise, one has applied Rule (ii) to obtain $M P^{\prime}$. But then again Rule (vii) is applicable on $M P^{\prime}=$ $\left(\square \gtrdot \ell_{0}, C, t,\left(\mathcal{M} \backslash\left\{f\left(t_{1}, \ldots, t_{m^{\prime}}\right) \gtrdot f\left(\ell_{1}^{\prime}, \ldots, \ell_{m^{\prime}}^{\prime}\right)\right\}\right) \cup\left\{t_{i} \gtrdot \ell_{i}^{\prime} \mid 1 \leq i \leq m^{\prime}\right\}\right.$, $\mu)$. The reason is again that the resulting matching problems that are considered in (vii) are equi-satisfiable as the only difference is between the one pair $f\left(t_{1}, \ldots, t_{m^{\prime}}\right) \gtrdot f\left(\ell_{1}^{\prime}, \ldots, \ell_{m^{\prime}}^{\prime}\right)$ and the set of pairs $\left\{s_{i} \gtrdot \ell_{i} \mid 1 \leq i \leq m\right\}$.

We are left to analyze all situations where $M P \Rightarrow M P_{1}$ and $M P \Rightarrow M P_{2}$ and $M P_{1} \neq M P_{2}$. If the corresponding rules that have been applied are (i) and (ii), or twice (ii) then two different decompositions have been performed to obtain $M P_{1}$ and $M P_{2}$. Then $M P_{1} \Rightarrow M P^{\prime}$ and $M P_{2} \Rightarrow M P^{\prime}$ for the matching problem $M P^{\prime}$ where both decompositions have been performed. If $M P_{1}$ is obtained by one of the decomposition Rules (i) or (ii) and $M P_{2}$ is obtained by instantiating the matching problem by (vi) then again one directly obtains a common successor $M P^{\prime}$ where both the decomposition and the instantiation have been performed. This is possible as decomposition and instantiation commute. If $M P_{1}$ is the result of applying Rule (ii) and $M P_{2}$ is the result of (viii) then again one can apply the corresponding other rule to achieve the same matching problem $M P^{\prime}$. This is due to the fact that decomposition yields equi-satisfiable matching problems which are checked by (viii). Finally, we have to consider the case where $M P_{1}=\left(\square \gtrdot \ell_{0}, C \mu, t \mu,\left\{s_{i} \mu \gtrdot \ell_{i} \mid 1 \leq i \leq m\right\}, \mu\right)$ is obtained by Rule (vi) and $M P_{2}=\left(C \gtrdot \ell_{0}, C \mu, t \mu,\left\{s_{i} \gtrdot \ell_{i} \mid 1 \leq i \leq m\right\}, \mu\right)$ is the result of (viii). Again, applying the corresponding other rule yields the same matching problem $\mathcal{M}^{\prime}=\left(C \mu \gtrdot \ell_{0}, C \mu \mu, t \mu \mu,\left\{s_{i} \mu \gtrdot \ell_{i} \mid 1 \leq i \leq m\right\}, \mu\right)$.

Using the above theorem allows us to transform any initial matching problem into $\perp$, $\top$, or into a matching problem in solved form. In the first two cases solvability is decided, but in the last case we still need a way to extract solvability. Note that these resulting matching problems are all of the form $M P=\left(D \gtrdot x_{0}, C, t,\left\{s_{1} \gtrdot x_{1}, \ldots, s_{m} \gtrdot x_{m}\right\}, \mu\right)$ where each $x_{i} \in \mathcal{V}$. Note that if all $x_{i}$ are different-which is always the case if one considers left-linear TRSs - then $M P$ is trivially solvable. One just can choose the solution $(n, k, \sigma)$ where $n=k=0$ and $\sigma=\left\{x_{0} / D[t], x_{1} / s_{1}, \ldots, x_{m} / s_{m}\right\}$.

The only problem arises if for some $i \neq j$ we have $x_{i}=x_{j}$. Then to choose $\sigma\left(x_{i}\right)=\sigma\left(x_{j}\right)$ one has to know that $s_{i} \mu^{k}=s_{j} \mu^{k}$ for some $k$. This so called identity problem already occurred in [15]. However, if $i=0$ then we have to answer a more difficult question, namely whether $D\left[t(C, \mu)^{n}\right] \mu^{k}=s_{j} \mu^{k}$. This new kind of problem is introduced as extended identity problem.
Definition 2 ((Extended) identity problems). An identity problem is a pair $\left(s \approx s^{\prime}, \mu\right)$. It is solvable iff there is some $k$ such that $s \mu^{k}=s^{\prime} \mu^{k}$.

An extended identity problem is a quadruple $(D \approx s, \mu, C, t)$. It is solvable iff there is a solution $(n, k)$ such that $D\left[t(C, \mu)^{n}\right] \mu^{k}=s \mu^{k}$.

We now can transform matching problems in solved form into an equivalent set of (extended) identity problems.
Theorem 2 (Transforming matching problems into identity problems). Let $M P=\left(D \gtrdot x, C, t,\left\{s_{1} \gtrdot x_{1}, \ldots, s_{m} \gtrdot x_{m}\right\}, \mu\right)$ be a matching problem in solved form. It is solvable iff each of the following identity problems is solvable.

- $\left(D \approx s_{i}, \mu, C, t\right)$ where $i$ is the least index such that $x=x_{i}$.
- $\left(s_{i} \approx s_{j}, \mu\right)$ for all $j$ where $i<j$ is the least index such that $x_{i}=x_{j}$.

Proof. First we prove the easy direction that if $M P$ is solvable then all identity problems are solvable. The reason is that any solution $(n, k, \sigma)$ of $M P$ satisfies the equations $D\left[t(C, \mu)^{n}\right] \mu^{k}=x \sigma$ and $s_{i} \mu^{k}=x_{i} \sigma$. Thus, whenever $x=x_{i}$ we obtain $D\left[t(C, \mu)^{n}\right] \mu^{k}=x \sigma=x_{i} \sigma=s_{i} \mu^{k}$ and hence, $(n, k)$ is a solution to ( $D \approx s_{i}, \mu, C, t$ ). Similarly, whenever $x_{i}=x_{j}$ then $s_{i} \mu^{k}=x_{i} \sigma=x_{j} \sigma=s_{j} \mu^{k}$ and hence, $\left(s_{i} \approx s_{j}, \mu\right)$ is solvable.

For the other direction suppose that all identity problems are solvable. We only consider the case that there is some least $i^{\prime}$ such that $x=x_{i^{\prime}}$. Hence, there is some $\left(n, k_{0}\right)$ such that $D\left[t(C, \mu)^{n}\right] \mu^{k_{0}}=s_{i} \mu^{k_{0}}$. Moreover, for all $i, j$ with $i<j$ and where $j$ is the least index such that $x_{i}=x_{j}$ we obtain some $k_{i j}$ such that $s_{i} \mu^{k_{i j}}=s_{j} \mu^{k_{i j}}$. We define $k$ to be the maximum of all these $k_{i j}$ 's and $k_{0}$. Moreover, we define $\sigma(x)=D\left[t(C, \mu)^{n}\right] \mu^{k}$ and $\sigma\left(x_{i}\right)=s_{i} \mu^{k}$. Then obviously, $(n, k, \sigma)$ is a solution for $M P$. The only missing part is to show that $\sigma$ is welldefined.

So, suppose there are $i, j$ such that $x_{i}=x_{j}$, but $s_{i} \mu^{k} \neq s_{j} \mu^{k}$. Note that for fixed $j$ there must be some least $i$ with these properties. Hence, for this $i$ we know that $s_{i} \mu^{k_{i j}}=s_{j} \mu^{k_{i j}}$. But since $k_{i j} \leq k$ we know that $s_{i} \mu^{k}=$ $s_{i} \mu^{k_{i j}} \mu^{k-k_{i j}}=s_{j} \mu^{k_{i j}} \mu^{k-k_{i j}}=s_{j} \mu^{k}$, a contradiction. Similarly, we obtain $\sigma(x)=$ $D\left[t(C, \mu)^{n}\right] \mu^{k}=D\left[t(C, \mu)^{n}\right] \mu^{k_{0}} \mu^{k-k_{0}}=s_{i} \mu^{k_{0}} \mu^{k-k_{0}}=s_{i} \mu^{k}=\sigma\left(x_{i}\right)$. Hence, $\sigma$ is indeed well-defined.

One might wonder why this theorem is sound as each solution of $\left(s_{i} \approx s_{j}, \mu\right)$ might yield a different $k_{i j}$. The key point is that the maximum of all these $k_{i j}$ 's is a solution for all identity problems $\left(s_{i} \approx s_{j}, \mu\right)$.

Note that [15] describes a decision procedure for solvability of identity problems $\left(s \approx s^{\prime}, \mu\right)$. Hence, we can already decide solvability of matching problems $(D \gtrdot x, C, t, \mathcal{M}, \mu)$ in solved form where the variable $x$ does not occur in $\mathcal{M}$. Nevertheless, for the general case we still need a technique to decide solvability of extended identity problems. Such a technique is described in the next section.

Example 2. Consider the TRS $\{\mathrm{f}(x) \rightarrow \mathrm{g}(\mathrm{g}(x, x), \mathrm{f}(\mathrm{s}(x))), \mathrm{g}(y, y) \rightarrow \mathrm{a}\}$ with the loop $t=\mathrm{f}(x) \rightarrow \mathrm{g}(\mathrm{g}(x, x), \mathrm{f}(\mathrm{s}(x)))=t(C, \mu)$ where $\mu=\{x / \mathrm{s}(x)\}$ and $C=$ $\mathrm{g}(\mathrm{g}(x, x), \square)$. One initial matching problem ( $\mathrm{g}(\mathrm{g}(x, x), \square) \gtrdot \mathrm{f}(x), C \mu, t \mu, \mu)$ is trivially not solvable due to a symbol clash. But the other initial matching problem $(\mathrm{g}(\mathrm{g}(x, x), \square) \gtrdot \mathrm{g}(y, y), C \mu, t \mu, \mu)$ is transformed into the matching problem $M P=(\square \gtrdot y, C \mu, t \mu,\{\mathrm{~g}(x, x) \gtrdot y\}, \mu)$. By Thm. 2 solvability of $M P$ is equivalent to solvability of the extended identity problem $(\square \approx \mathrm{g}(x, x), \mu, C \mu, t \mu)$.

## 5 Deciding Solvability of Extended Identity Problems

In this section we describe a decision procedure for solvability of extended identity problems. To this end we first introduce the notion of a trace.

Definition 1 (Traces). The trace of term $t$ w.r.t. position $p$ is the sequence of function symbols and indices that are passed when moving from $\varepsilon$ to $p$ in $t$ :

$$
\operatorname{trace}(p, t)= \begin{cases}\varepsilon & \text { if } t=x \text { or } p=\varepsilon \\ f i \operatorname{trace}\left(q, t_{i}\right) & \text { if } p=i q \text { and } t=f\left(t_{1}, \ldots, t_{n}\right)\end{cases}
$$

The trace of a context $C$ with $\left.C\right|_{p}=\square$ is $\operatorname{trace}(C)=\operatorname{trace}(p, C)$. The set of all traces of a term is $\operatorname{Traces}(t)=\{\operatorname{trace}(p, t) \mid p \in \mathcal{P} \operatorname{os}(t)\}$.
Lemma 1 (Properties of traces).
(i) $\operatorname{trace}(p q, C[t])=\operatorname{trace}(C) \operatorname{trace}(q, t)$ if $\left.C\right|_{p}=$
(ii) $\operatorname{trace}(p, t)=\operatorname{trace}(p, t \mu)$ if $p \in \mathcal{P} \operatorname{os}(t)$
(iii) $\operatorname{trace}\left(p^{n} q, t(C, \mu)^{n}\right)=\operatorname{trace}(C)^{n} \operatorname{trace}\left(q, t \mu^{n}\right)$ if $\left.C\right|_{p}=\square$ and $q \in \mathcal{P} \operatorname{os}(t)$
(iv) $\operatorname{trace}(p, t) \in \operatorname{Traces}(t)$ if $\operatorname{trace}(p q, t) \in \operatorname{Traces}(t)$

Proof. (i) We perform induction over $C$. For the base case we have $C=$ and thus $p=\varepsilon$. Hence, $\operatorname{trace}(p q, C[t])=\operatorname{trace}(q, t)=\operatorname{trace}(C) \operatorname{trace}(q, t)$. In the step case we have $C=f(\ldots, D, \ldots)$ for some context $D$ with $\left.C\right|_{i}=$ $D,\left.D\right|_{p^{\prime}}=\square$, and $\left.C\right|_{p}=\square$. Hence $p=i p^{\prime}$. The IH is $\operatorname{trace}\left(p^{\prime} q, D[t]\right)=$ $\operatorname{trace}(D) \operatorname{trace}(q, t)$. We conclude

$$
\begin{align*}
\operatorname{trace}(p q, C[t]) & =\operatorname{trace}\left(i p^{\prime} q, f(\ldots, D[t], \ldots)\right) \\
& =f i \operatorname{trace}\left(p^{\prime} q, D[t]\right) \\
& =f i \operatorname{trace}(D) \operatorname{trace}(q, t)  \tag{byIH}\\
& =\operatorname{trace}(C) \operatorname{trace}(q, t)
\end{align*}
$$

(ii) We use induction over $p$. For the base case we have $p=\varepsilon$. Thus $\operatorname{trace}(p, t)=$ $\varepsilon=\operatorname{trace}(p, t \mu)$. In the step case we have $p=i q$ for some position $q$. From $p \in \mathcal{P o s}(t)$ we conclude that $t=f\left(t_{1}, \ldots, t_{n}\right)$ where $q \in \mathcal{P o s}\left(t_{i}\right)$. The IH is $\operatorname{trace}\left(q, t_{i}\right)=\operatorname{trace}\left(q, t_{i} \mu\right)$. We finish this part by

$$
\begin{aligned}
\operatorname{trace}(p, t) & =\operatorname{trace}\left(i q, f\left(t_{1}, \ldots, t_{n}\right)\right) \\
& =f i \operatorname{trace}\left(q, t_{i}\right) \\
& =f i \operatorname{trace}\left(q, t_{i} \mu\right) \\
& =\operatorname{trace}\left(i q, f\left(t_{1} \mu, \ldots, t_{n} \mu\right)\right) \\
& =\operatorname{trace}\left(i q, f\left(t_{1}, \ldots, t_{n}\right) \mu\right) \\
& =\operatorname{trace}(p, t \mu)
\end{aligned}
$$

(iii) We use induction over $n$. For the base case we have $n=0$ and hence $\operatorname{trace}\left(p^{n} q, t(C, \mu)^{n}\right)=\operatorname{trace}(q, t)=\operatorname{trace}(C)^{n} \operatorname{trace}\left(q, t \mu^{n}\right)$. In the step case we have $n=k+1,\left.C\right|_{p}=\square$, and $q \in \mathcal{P} \operatorname{os}(t)$; and hence also $q \in \mathcal{P o s}\left(t \mu^{k}\right)$. By an easy inductive argument it can also be shown that this implies $p^{k} q \in$ $\mathcal{P o s}\left(t(C, \mu)^{k}\right)$. The IH is $\operatorname{trace}\left(p^{k} q, t(C, \mu)^{k}\right)=\operatorname{trace}(C)^{k} \operatorname{trace}\left(q, t \mu^{k}\right)$. Thus,

$$
\begin{array}{rlr}
\operatorname{trace}(C)^{k+1} \operatorname{trace}\left(q, t \mu^{k+1}\right) & =\operatorname{trace}(C) \operatorname{trace}(C)^{k} \operatorname{trace}\left(q, t \mu^{k+1}\right) \\
& =\operatorname{trace}(C) \operatorname{trace}(C)^{k} \operatorname{trace}\left(q,\left(t \mu^{k}\right) \mu\right) \\
& =\operatorname{trace}(C) \operatorname{trace}(C)^{k} \operatorname{trace}\left(q, t \mu^{k}\right) & \quad(\text { by (ii)) } \\
& =\operatorname{trace}(C) \operatorname{trace}\left(p^{k} q, t(C, \mu)^{k}\right) & \quad(\text { by IH) } \\
& =\operatorname{trace}(C) \operatorname{trace}\left(p^{k} q,\left(t(C, \mu)^{k}\right) \mu\right) & \quad(\text { by (ii)) } \\
& =\operatorname{trace}\left(p p^{k} q, C\left[\left(t(C, \mu)^{k}\right) \mu\right]\right) & \quad \text { (by (i)) } \\
& =\operatorname{trace}\left(p^{k+1} q, t(C, \mu)^{k+1}\right) . &
\end{array}
$$

(iv) From $\operatorname{trace}(p q, t) \in \operatorname{Traces}(t)$ we know that $p q \in \mathcal{P o s}(t)$ and hence also $p \in \mathcal{P o s}(t)$. Thus trace $(p, t) \in \operatorname{Traces}(t)$ by the definition of $\operatorname{Traces}(t)$.

In the following algorithm to decide solvability of extended identity problems, a Boolean disjunction over non-extended identity problems represents solvability of at least one of these identity problems.

Definition 2 (Decision procedure for extended identity problems). Let $(D \approx s, \mu, C, t)$ be an extended identity problem where $\left.D\right|_{q}=\square,\left.C\right|_{p}=\square$.
(i) if $\operatorname{trace}(D) \operatorname{trace}(C)^{*} \nsubseteq \bigcup_{i \in \mathbb{N}} \operatorname{Traces}\left(s \mu^{i}\right)=$ : $S$ then there is some $m$ such that $\operatorname{trace}(D) \operatorname{trace}(C)^{m} \notin S$; return $\bigvee_{n<m}\left(D\left[t(C, \mu)^{n}\right] \approx s, \mu\right)$
(ii) if $\operatorname{trace}(C)^{*} \nsubseteq \bigcup_{i \in \mathbb{N}} \operatorname{Traces}\left(t \mu^{i}\right)$ then return "not solvable"
(iii) let $x$ be a variable which infinitely often occurs in $\left.s\right|_{p_{0}},\left.s \mu^{1}\right|_{p_{1}},\left.s \mu^{2}\right|_{p_{2}}, \ldots$ where each $p_{i}$ is that prefix of $q p^{\omega}$ which satisfies $p_{i} \in \mathcal{P} \operatorname{os}\left(s \mu^{i}\right)$; let $i$ be the minimal number such that $\left.s \mu^{i}\right|_{p_{i}}=x$
(iv) let $j$ be minimal such that $\left.t \mu^{j}\right|_{q_{j}}=x$ where $q_{j}$ is that prefix of $p^{\omega}$ which satisfies $q_{j} \in \operatorname{Pos}\left(t \mu^{j}\right)$; if there is no such $j$ then return "not solvable"
(v) return $\bigvee_{n \leq \max \left(\frac{\left|p_{i}\right|-\left|q q_{j}\right|}{|p|}, j\right)}\left(D\left[t(C, \mu)^{n}\right] \approx s, \mu\right)$

We will explain the algorithm in detail within the proof of the following theorem. Afterwards, we present algorithms to automate the non-trivial steps.

Theorem 1. The algorithm of Def. 2 is sound and terminates.
Proof. Termination of the algorithm is obvious. We only remark that the disjunction in Step (v) is finite, since $|p|>0$ by the assumption $C \neq \square$.

To show soundness of the algorithm first recall the definition of solvability of $(D \approx s, \mu, C, t)$. This extended identity problem is solvable iff there is a solution $(n, k)$ such that $D\left[t(C, \mu)^{n}\right] \mu^{k}=s \mu^{k}$. We observe two properties: first, whenever $(n, k)$ is a solution then $\left(n, k+k^{\prime}\right)$ is also a solution. And second, if we fix $n$ then the extended identity problem is solvable iff the identity problem $\left(D\left[t(C, \mu)^{n}\right] \approx s, \mu\right)$ is solvable. From the second observation we conclude that if one can bound the value of $n$, then one can reduce solvability of extended identity problems to solvability of identity problems and is done. And computing these bounds on $n$ is basically all the algorithm does (in Steps (i) and (v)).

The first idea to extract a bound is to consider how the term $D\left[t(C, \mu)^{n}\right] \mu^{k}$ grows if $n$ is increased. Looking at Fig. 1 on page 3 or using Lemma 1 we see that $D\left[t(C, \mu)^{n}\right] \mu^{k}$ has the trace trace $(D) \operatorname{trace}(C)^{n}$. Thus, if $(n, k)$ is a solution then $s \mu^{k}$ must have the same trace. Hence, if trace $(D) \operatorname{trace}(C)^{m} \notin \bigcup_{i \in \mathbb{N}} \operatorname{Traces}\left(s \mu^{i}\right)=$ $S$ then $n<m$. This proves soundness of Step (i).

So, after Step (i) we can assume trace $(D) \operatorname{trace}(C)^{*} \subseteq S$. Hence, if we increase the $k$ of $s \mu^{k}$ then this term grows along the (infinite) trace trace $(D) \operatorname{trace}(C)^{\omega}$. Using the first observation we know that for every solution of the extended identity problem we can increase $k$ arbitrarily. Thus, $D\left[t(C, \mu)^{n}\right] \mu^{k}$ (which is the same term as $s \mu^{k}$ ) also has to contain longer and longer parts of the trace trace $(D) \operatorname{trace}(C)^{\omega}$ when increasing $k$. Hence, whenever the extended identity problem is solvable then trace $(C)^{*} \subseteq \bigcup_{i \in \mathbb{N}} \operatorname{Traces}\left(t \mu^{i}\right)=: T$ which shows soundness of Step (ii).

If the decision procedure arrives at Step (iii) then both trace $(D) \operatorname{trace}(C)^{*} \subseteq S$ and trace $(C)^{*} \subseteq T$, i.e., when increasing $k$ we see no difference of function symbols of the terms $D\left[t(C, \mu)^{n}\right] \mu^{k}$ and $s \mu^{k}$ along the path $q p^{\omega}$. However, there still might be a difference between $D\left[t(C, \mu)^{n}\right] \mu^{k}$ and $s \mu^{k}$ for each finite value $k$. For example, different variables may be used to increase the terms along the path $q p^{\omega}$ as in $D=\square, C=f(\square), t=x, s=y, \mu=\{x / f(x), y / f(y)\}$. Or one of the terms is always a bit larger as the other one as in $D=\square, C=f(\square), t=$ $f(x), s=x, \mu=\{x / f(x)\}$. To detect the situation of different variables, Step (iv) is used, and the latter situation is done via Step (v).

As we are interested in the variables in Steps (iv) and (v), we first compute a variable $x$ in Step (iii) which infinitely often occurs along the path $q p^{\omega}$ in the terms $s, s \mu, s \mu^{2}, \ldots$. This variable must exist, since $\mu$ has finite domain and trace $(D)$ trace $(C)^{*} \subseteq S$. This immediately proves soundness of Step (iv) since
whenever $D\left[t(C, \mu)^{n}\right] \mu^{k}=s \mu^{k}$ then by the first observation we can choose $k$ high enough such that $x=\left.s \mu^{k}\right|_{p_{k}}=\left.D\left[t(C, \mu)^{n}\right] \mu^{k}\right|_{p_{k}}$ where $p_{k} \leq q p^{\omega}$, i.e., $p_{k}$ must be of the form $q p^{n} q^{\prime}$ where $q^{\prime}$ is a prefix of $p^{\omega}$. Thus, $x=\left.D\left[t(C, \mu)^{n}\right] \mu^{k}\right|_{q p^{n} q^{\prime}}=$ $\left.t(C, \mu)^{n} \mu^{k}\right|_{p^{n} q^{\prime}}=\left.t \mu^{n+k}\right|_{q^{\prime}}$ shows that Step (iv) cannot stop the algorithm with "not solvable".

The main idea of Step (v) is as in Step (i) to bound $n$ but now for a different reason. Observe that whenever we increase $n$ then $s \mu^{k}$ stays the same whereas $D\left[t(C, \mu)^{n}\right] \mu^{k}$ has the subterm $t \mu^{n+k}$ which depends on $n$. And since trace $(C)^{*} \subseteq$ $T$ we know that with larger $n$ also the terms $t \mu^{n+k}$ become larger. Thus, there must be a limit where the size of $s \mu^{k}$ is reached and it is of no use to search for larger values of $n$. And this limit turns out to be $m:=\max \left(\frac{\left|p_{i}\right|-\left|q q_{j}\right|}{|p|}, j\right)$ which proves soundness of Step (v).

Proof. We only present those details that have been omitted in the proof of Thm. 1 on page 16.

To prove soundness of Step (ii) we can already assume trace $(D) \operatorname{trace}(C)^{*} \subseteq$ $S$ due to $\operatorname{Step}$ (i). We have to show that trace $(C)^{*} \subseteq \bigcup_{i \in \mathbb{N}} \operatorname{Traces}\left(t \mu^{i}\right)=: T$ whenever the extended identity problem is solvable. So, suppose there is some $m$ such that trace $(C)^{m} \notin T$ and that $D\left[t(C, \mu)^{n}\right] \mu^{k}=s \mu^{k}$. Using the latter equality in combination with the first observation and trace $(D) \operatorname{trace}(C)^{*} \subseteq S$, we can assume $k$ to be large enough such that

$$
\operatorname{trace}(D) \operatorname{trace}(C)^{n+m} \in \operatorname{Traces}\left(s \mu^{k}\right)=\operatorname{Traces}\left(D\left[t(C, \mu)^{n}\right] \mu^{k}\right)
$$

Recall that $\left.D\right|_{q}=\square$ and $\left.C\right|_{p}=\square$. Using Lemma 1 we derive

$$
\begin{aligned}
\operatorname{trace}(D) \operatorname{trace}(C)^{n+m} & =\operatorname{trace}\left(q p^{n+m}, D\left[t(C, \mu)^{n}\right] \mu^{k}\right) \\
& =\operatorname{trace}\left(D \mu^{k}\right) \operatorname{trace}\left(p^{n+m}, t(C, \mu)^{n} \mu^{k}\right) \\
& =\operatorname{trace}\left(D \mu^{k}\right) \operatorname{trace}\left(C \mu^{k}\right)^{n} \operatorname{trace}\left(p^{m}, t \mu^{n+k}\right) \\
& =\operatorname{trace}(D) \operatorname{trace}(C)^{n} \operatorname{trace}\left(p^{m}, t \mu^{n+k}\right)
\end{aligned}
$$

Thus, trace $(C)^{m}=\operatorname{trace}\left(p^{m}, t \mu^{n+k}\right)$ must hold in contradiction to the assumption trace $(C)^{m} \notin T$. Hence, Step (ii) is sound.

To show that $m=\max \left(\frac{\left|p_{i}\right|-\left|q q_{j}\right|}{|p|}, j\right)$ really is a limit on $n$ in Step (v), we assume that there is a solution $(n, k)$ where $n>m$. Note that when executing Step (v) of the decision procedure, one has obtained minimal values $i$ and $j$ such that $\left.s \mu^{i}\right|_{p_{i}}=\left.t \mu^{j}\right|_{q_{j}}=x$ for $p_{i}<q p^{\omega}$ and $q_{j}<p^{\omega}$. Due to the first observation we may assume $k \geq i$. Then $D\left[t(C, \mu)^{n}\right] \mu^{k}=s \mu^{k}, k=i+i^{\prime}$ for some $i^{\prime} \geq 0$, $n=j+j^{\prime}$ for some $j^{\prime} \geq 0, n>\frac{\left|p_{i}\right|-\left|q q_{j}\right|}{|p|}$, and hence, $\left|q p^{n} q_{j}\right|>\left|p_{i}\right|$. Since both $q p^{n} q_{j}$ and $p_{i}$ are prefixes of $q p^{\infty}$ we know that there is some $q^{\prime} \neq \varepsilon$ such that
$q p^{n} q_{j}=p_{i} q^{\prime}$. Using all these equations we conclude

$$
\begin{aligned}
x \mu^{i^{\prime}} \mu^{i+j^{\prime}} & =x \mu^{j^{\prime}+k} \\
& =\left.t \mu^{j}\right|_{q_{j}} \mu^{j^{\prime}+k} \\
& =\left.t \mu^{k} \mu^{n}\right|_{q_{j}} \\
& =\left.t \mu^{k}\left(C \mu^{k}, \mu\right)^{n}\right|_{p^{n} q_{j}} \\
& =\left.t(C, \mu)^{n} \mu^{k}\right|_{p^{n} q_{j}} \\
& =\left.D\left[t(C, \mu)^{n}\right] \mu^{k}\right|_{q p^{n} q_{j}} \\
& =\left.s \mu^{k}\right|_{q p^{n} q_{j}} \\
& =\left.s \mu^{i+i^{\prime}}\right|_{p_{i} q^{\prime}} \\
& =\left.x \mu^{i^{\prime}}\right|_{q^{\prime}}
\end{aligned}
$$

which gives rise to a contradiction: the instance $x \mu^{i^{\prime}} \mu^{i+j^{\prime}}$ of the term $x \mu^{i^{\prime}}$ cannot be identical to the proper subterm $\left.x \mu^{i^{\prime}}\right|_{q^{\prime}}$ of $x \mu^{i^{\prime}}$.

Input: $D, C, s, \mu$
Output: minimal $m$ such that trace $(D) \operatorname{trace}(C)^{m} \notin \bigcup_{i \in \mathbb{N}} \operatorname{Traces}\left(s \mu^{i}\right)$ or $\infty$, otherwise
(1) $m:=0, E:=D, t:=s, S:=\varnothing$
(2) if $E=f\left(\ldots E^{\prime} \ldots\right)$ and $t=f(\ldots)$ where $\left.E\right|_{i}=E^{\prime}$ then $E:=E^{\prime}, t:=\left.t\right|_{i}$, goto (2)
(3) if $E=g(\ldots), t=f(\ldots)$, and $g \neq f$ then return $m$
(4) if $E \neq \square$ and $t=x \notin \mathcal{V}_{\text {incr }}(\mu)$ then return $m$
(5) if $E \neq \square$ and $t=x \in \mathcal{V}_{\text {incr }}(\mu)$ then $t:=t \mu$, goto (2)
(6) if $E=\square$ and $t \in S$ then return $\infty$
(7) if $E=\square$ and $t \notin S$ then $S:=S \cup\{t\}, m:=m+1, E:=C$, goto (2)

Fig. 2. $\operatorname{clash}(D, C, s, \mu)$.

For the automation one can use the algorithm of [15] for solvability of nonextended identity problems. To check whether trace $(D) \operatorname{trace}(C)^{*} \subseteq S$ in Step (i) one can check whether $\operatorname{clash}(D, C, s, \mu)=\infty$ (cf. Fig. 2) which also delivers the required number $m$ in case that trace $(D)$ trace $(C)^{*} \nsubseteq S$. Of course, clash can also be used for the test in Step (ii) where we call clash $(\square, C, t, \mu)$.
Theorem 2. The algorithm clash terminates and is sound.
Proof. Termination of the algorithm in Fig. 2 can be proven as follows. First note that the term $t$ will always be a subterm of $s$ or a subterm of $x \mu$ where $x \in \operatorname{Dom}(\mu)$. Hence, it is not possible to perform Step (7) infinitely often. Thus, after the last application of (7) the context $E$ can only become smaller. Hence, Step (2) cannot be performed infinitely often, too. And since every application of Step (5) eventually triggers one of the Steps (2) or (3), termination is proven.

Soundness is proven with the help of the following invariant before Step (2).

There is some position $p$ and number $k$ such that $\operatorname{trace}(D) \operatorname{trace}(C)^{m}=\operatorname{trace}\left(p, s \mu^{k}\right) \operatorname{trace}(E)$ and $t=\left.s \mu^{k}\right|_{p}$.

First note that the invariant is satisfied initially by taking $p=\varepsilon$ and $k=0$.
In the remainder of this proof (and also in the upcoming proofs), $\bar{v}$ is used for some value $v$ before applying a step and $\underline{v}$ for the same value after applying the step. As induction hypothesis we always have some $\bar{p}$ and $\bar{k}$ such that $\operatorname{trace}(D) \operatorname{trace}(C)^{m}=\operatorname{trace}\left(\bar{p}, s \mu^{\bar{k}}\right) \operatorname{trace}(\bar{E})$, and $\bar{t}=\left.s \sigma^{\bar{k}}\right|_{\bar{p}}$.

For Step (2) we know that $\bar{E}=f(\ldots), \underline{E}=\left.\bar{E}\right|_{i}, \bar{t}=f(\ldots)$, and $\underline{t}=\left.\bar{t}\right|_{i}$. Hence,

$$
\begin{aligned}
\operatorname{trace}(D) \operatorname{trace}(C)^{m} & =\operatorname{trace}\left(\bar{p}, s \mu^{\bar{k}}\right) \operatorname{trace}(\bar{E}) \\
& =\operatorname{trace}\left(\bar{p}, s \mu^{\bar{k}}\right) f i \operatorname{trace}(\underline{E}) \\
& =\operatorname{trace}\left(\bar{p}, s \mu^{\bar{k}}\right) \operatorname{trace}(i, \bar{t}) \operatorname{trace}(\underline{E}) \\
& =\operatorname{trace}\left(\bar{p}, s \mu^{\bar{k}}\right) \operatorname{trace}\left(i, s \sigma^{\bar{k}} \mid \bar{p}\right) \operatorname{trace}(\underline{E}) \\
& =\operatorname{trace}\left(\bar{p} i, s \mu^{\bar{k}}\right) \operatorname{trace}(\underline{E}) \\
& =\operatorname{trace}\left(\underline{p}, s \mu^{\underline{k}}\right) \operatorname{trace}(\underline{E})
\end{aligned}
$$

if we define $\underline{p}=\bar{p} i$ and $\underline{k}=\bar{k}$. Moreover, we also obtain $\underline{t}=\left.\bar{t}\right|_{i}=\left.\left.s \sigma^{\bar{k}}\right|_{\bar{p}}\right|_{i}=\left.s \sigma^{\underline{k}}\right|_{\underline{p}}$.
For Step (7) we know that $\underline{E}=C$ and $\bar{E}=\square$, i.e., $\operatorname{trace}(D) \operatorname{trace}(C)^{m}=$ $\operatorname{trace}\left(\bar{p}, s \mu^{\bar{k}}\right) \operatorname{trace}(\bar{E})=\operatorname{trace}\left(\bar{p}, s \mu^{\bar{k}}\right)$. Choosing $\underline{k}=\bar{k}$ and $\underline{p}=\bar{p}$ yields

$$
\operatorname{trace}(D) \operatorname{trace}(C)^{m+1}=\operatorname{trace}\left(\bar{p}, s \mu^{\bar{k}}\right) \operatorname{trace}(C)=\operatorname{trace}\left(\underline{p}, s \mu^{\underline{k}}\right) \operatorname{trace}(\underline{E})
$$

Moreover, all other properties of the invariant are trivially satisfied.
In Step (5) where $\underline{t}=\bar{t} \mu$ we choose $\underline{k}=\bar{k}+1$ and $\underline{p}=\bar{p}$. Then, $\underline{t}=\bar{t} \mu=$ $\left.s \mu^{\bar{k}}\right|_{\bar{p}} \mu=\left.s \mu^{\bar{k}+1}\right|_{\bar{p}}=\left.s \mu^{\underline{k}}\right|_{\underline{p}}$ And since $\bar{p} \in \operatorname{Pos}\left(s \mu^{\bar{k}}\right)$ we know that $\operatorname{trace}\left(\bar{p}, s \mu^{\bar{k}}\right)=$ $\operatorname{trace}\left(\bar{p}, s \mu^{\bar{k}+1}\right)=\operatorname{trace}\left(\underline{p}, s \mu^{\underline{k}}\right)$. Thus, $\operatorname{trace}(D) \operatorname{trace}(C)^{m}=\operatorname{trace}\left(\underline{p}, s \mu^{\underline{k}}\right) \operatorname{trace}(\underline{E})$ is immediately obtained.

Soundness of Step (3) is due to the fact that

$$
\operatorname{trace}(D) \operatorname{trace}(C)^{m}=\operatorname{trace}\left(\bar{p}, s \mu^{\bar{k}}\right) \operatorname{trace}(\bar{E})=\operatorname{trace}\left(\bar{p}, s \mu^{\bar{k}}\right) g \ldots
$$

whereas

$$
\operatorname{trace}\left(\bar{p} j, s \mu^{k}\right)=\operatorname{trace}\left(\bar{p}, s \mu^{\bar{k}}\right) f \ldots
$$

for all $1 \leq j \leq \operatorname{ar}(f)$. Thus, trace $(D) \operatorname{trace}(C)^{m} \notin \bigcup_{i \in \mathbb{N}} \operatorname{Traces}\left(s \mu^{i}\right)$. Moreover, if $m=0$ then minimality of $m$ is obvious. For the other case where $m>0$ one can easily see that $\bar{E}$ must be a sub-context of $C$. Hence, trace $(D) \operatorname{trace}(C)^{m}=$ $\operatorname{trace}\left(\bar{p}, s \mu^{\bar{k}}\right) \operatorname{trace}(\bar{E})$ implies that $\operatorname{trace}(D) \operatorname{trace}(C)^{m-1} \leq \operatorname{trace}\left(\bar{p}, s \mu^{\bar{k}}\right)$. Thus, $\operatorname{trace}(D) \operatorname{trace}(C)^{m-1} \in \operatorname{Traces}\left(s \mu^{\bar{k}}\right)$ shows minimality of $m$.

Soundness of Step (4) is proven in the similar way as (3). Here, the problem is that $\operatorname{trace}\left(\bar{p}, s \mu^{\bar{k}}\right)$ is a maximal trace within $\bigcup_{i \in \mathbb{N}} \operatorname{Traces}\left(s \mu^{i}\right)$ since $\left.s \mu^{k^{\prime}}\right|_{\bar{p}}$ is a variable for all $k^{\prime} \geq \bar{k}$. Hence, it cannot be extended by $f \ldots$ which would be required since $\operatorname{trace}(\bar{E})=f \ldots$ as $\bar{E}=f(\ldots) \neq \square$.

Finally, for Step (6) we obtain $\operatorname{trace}(D) \operatorname{trace}(C)^{m}=\operatorname{trace}\left(\bar{p}, s \mu^{\bar{k}}\right),\left.s \mu^{\bar{k}}\right|_{\bar{p}}=t$ as in Step (7). And moreover, since $\bar{t} \in S$ there must be some $m^{\prime}<m, k^{\prime}$, and $p^{\prime}$ such that $\operatorname{trace}(D) \operatorname{trace}(C)^{m^{\prime}}=\operatorname{trace}\left(p^{\prime}, s \mu^{k^{\prime}}\right)$ and $\left.s \mu^{k^{\prime}}\right|_{p^{\prime}}=t$. Let $q, q^{\prime}$ such that $\left.C\right|_{q}=\left.D\right|_{q^{\prime}}=\square$. Thus, $\bar{p}=q^{\prime} q^{m}, p^{\prime}=q^{\prime} q^{m^{\prime}}$, and $m=m^{\prime}+m^{\prime \prime}$ for some $m^{\prime \prime}>0$. Since

$$
\left.s \mu^{k^{\prime}}\right|_{q^{\prime} q^{m^{\prime}}}=\left.s \mu^{k^{\prime}}\right|_{p^{\prime}}=t=\left.s \mu^{\bar{k}}\right|_{\bar{p}}=\left.s \mu^{\bar{k}}\right|_{q^{\prime} q^{m^{\prime}+m^{\prime \prime}}}
$$

we additionally know that $\bar{k}=k^{\prime}+k^{\prime \prime}$ for some $k^{\prime \prime}>0$. We conclude

$$
\begin{aligned}
\operatorname{trace}(D) \operatorname{trace}(C)^{m^{\prime}} \operatorname{trace}(C)^{m^{\prime \prime}} & =\operatorname{trace}(D) \operatorname{trace}(C)^{m} \\
& =\operatorname{trace}\left(q^{\prime} q^{m}, s \mu^{\bar{k}}\right) \\
& =\operatorname{trace}\left(q^{\prime} q^{m^{\prime}+m^{\prime \prime}}, s \mu^{k^{\prime}+k^{\prime \prime}}\right) \\
& =\operatorname{trace}\left(q^{\prime} q^{m^{\prime}}, s \mu^{k^{\prime}+k^{\prime \prime}}\right) \operatorname{trace}\left(q^{m^{\prime \prime}},\left.s \mu^{k^{\prime}+k^{\prime \prime}}\right|_{q^{\prime} q^{m^{\prime}}}\right) \\
& =\operatorname{trace}\left(q^{\prime} q^{m^{\prime}}, s \mu^{k^{\prime}}\right) \operatorname{trace}\left(q^{m^{\prime \prime}},\left.s \mu^{k^{\prime}}\right|_{q^{\prime} q^{m^{\prime}} \mu^{k^{\prime \prime}}}\right) \\
& =\operatorname{trace}(D) \operatorname{trace}(C)^{m^{\prime}} \operatorname{trace}\left(q^{m^{\prime \prime}}, t \mu^{k^{\prime \prime}}\right)
\end{aligned}
$$

and hence, $\operatorname{trace}(C)^{m^{\prime \prime}}=\operatorname{trace}\left(q^{m^{\prime \prime}}, t \mu^{k^{\prime \prime}}\right)$. By instantiating $\mu$ for $k^{\prime \prime}$ more times, we can thus append an additional trace $(C)^{m^{\prime \prime}}$ to the trace:

$$
\begin{aligned}
\operatorname{trace}\left(\bar{p} q^{m^{\prime \prime}}, s \mu^{\bar{k}+k^{\prime \prime}}\right) & =\operatorname{trace}\left(\bar{p}, s \mu^{\bar{k}+k^{\prime \prime}}\right) \operatorname{trace}\left(q^{m^{\prime \prime}},\left.s \mu^{\bar{k}+k^{\prime \prime}}\right|_{q^{\prime} q^{m}}\right) \\
& =\operatorname{trace}\left(\bar{p}, s \mu^{\bar{k}}\right) \operatorname{trace}\left(q^{m^{\prime \prime}}, t \mu^{k^{\prime \prime}}\right) \\
& =\operatorname{trace}(D) \operatorname{trace}(C)^{m} \operatorname{trace}\left(q^{m^{\prime \prime}}, t \mu^{k^{\prime \prime}}\right) \\
& =\operatorname{trace}(D) \operatorname{trace}(C)^{m} \operatorname{trace}(C)^{m^{\prime \prime}}
\end{aligned}
$$

Iterating in this way yields trace $(D) \operatorname{trace}(C)^{m+j m^{\prime \prime}} \in \bigcup_{i \in \mathbb{N}} \operatorname{Traces}\left(s \mu^{i}\right)=: T$ for all $j \in \mathbb{N}$. And this implies trace $(D) \operatorname{trace}(C)^{*} \subseteq T$ since whenever a trace is in $T$ then also all its prefixes are contained in $T$.

For Step (iii) of the decision procedure the function $\operatorname{var}_{\infty}(s, q, p, \mu)$ in Fig. 3 computes a triple $\left(x, p_{i}, i\right)$ such that $x$ infinitely often occurs in the sequence $\left.s\right|_{p_{0}},\left.s \mu\right|_{p_{1}},\left.s \mu^{2}\right|_{p_{2}}, \ldots$ where each $p_{k} \in \mathcal{P o s}\left(s \mu^{k}\right)$ is the maximal prefix of $q p^{\infty}$ and $i$ is the smallest number such that $\left.s \mu^{i}\right|_{p_{i}}=x$. Note that the precondition is satisfied, since $\operatorname{var}_{\infty}$ is only called if trace $(D) \operatorname{trace}(C)^{*} \subseteq \bigcup_{i \in \mathbb{N}} \operatorname{Traces}\left(s \mu^{i}\right)$ which directly implies $q p^{*} \subseteq \bigcup_{i \in \mathbb{N}} \mathcal{P} \operatorname{os}\left(s \mu^{i}\right)$.

Preconditions: $p \neq \varepsilon$ and $q p^{*} \subseteq \bigcup_{j \in \mathbb{N}} \mathcal{P o s}\left(s \mu^{j}\right)$
Input: $s, q, p, \mu$
Output: $\left(x, p_{i}, i\right)$ such that $x$ infinitely often occurs in terms $s \mu^{j}$ along path $q p^{\omega}$ where $i$ and $p_{i}<q p^{\omega}$ are minimal such that $\left.s \mu^{i}\right|_{p_{i}}=x$
(1) $i:=0, u:=s, p_{\text {start }}:=\varepsilon, p_{\text {end }}:=q, S:=\varnothing$
(2) if $u=x$ and $\left(x,-, p_{\text {end }},-\right) \in S$ then return $\left(x, p_{\text {start }}{ }^{\prime}, i^{\prime}\right)$ where $i^{\prime}$ is the smallest value such that $\left(x, p_{\text {start }}{ }^{\prime},-, i^{\prime}\right) \in S$
(3) if $u=x$ and $\left(x,-, p_{\text {end }},-\right) \notin S$ then $u:=u \mu, S:=S \cup\left\{\left(x, p_{\text {start }}, p_{\text {end }}, i\right)\right\}, i:=i+1$, goto (2)
(4) (a) if $p_{\text {end }}=\varepsilon$ then $p_{\text {end }}:=p$
(b) let $p_{\text {end }}=j p^{\prime}$ and $u=f\left(u_{1}, \ldots, u_{n}\right) ; u:=u_{j}, p_{\text {end }}:=p^{\prime}, p_{\text {start }}:=p_{\text {start }} j$, goto (2)

Fig. 3. $\operatorname{var}_{\infty}(s, q, p, \mu)$

Theorem 3. The algorithm $\mathrm{var}_{\infty}$ terminates and is sound.

Proof. Using the preconditions that $q p^{*} \subseteq \bigcup_{i \in \mathbb{N}} \mathcal{P} \operatorname{os}\left(s \mu^{i}\right)$ and $p \neq \varepsilon$, termination of the algorithm in Fig. 3 can be shown as follows. First note that $u$ will always be a subterm of $s$ or of $x \mu$ for some $x \in \operatorname{Dom}(\mu)$. Moreover, $p_{\text {end }}$ will always be a suffix of $q$ or of $p$. As long as $u$ is not a variable, Step (4-b) makes sure that at some point a variable is reached (using the preconditions). Further note that Step (3) can only be used finitely often since $\mu$ has finite domain and since there are only finitely many possible values for $p_{\text {end }}$. Hence, eventually Step (2) is applied and the algorithm stops.

To prove soundness we first introduce the notion of a minimal extension: the minimal extension of position $q^{\prime}$ w.r.t. the positions $q$ and $p$ is defined if $q^{\prime}$ is a prefix of $q p^{\omega}$; then there is a minimal $k$ such that $q^{\prime}$ is a prefix of $q p^{k}$ and the minimal extension of $q^{\prime}$ is the unique position $p^{\prime}$ such that $q^{\prime} p^{\prime}=q p^{k}$. Since $p$ and $q$ are fixed within this proof we just speak of the minimal extension of $q^{\prime}$ and omit the "w.r.t. $q$ and $p$ ".

We prove the following three invariants which are established before Step (2).
(i) $p_{\text {end }}$ is the minimal extension of $p_{\text {start }}$.
(ii) It holds that $u=\left.s \mu^{i}\right|_{p_{\text {start }}}$.
(iii) For all $j<i$ it holds that $\left(x, p_{\text {start }}{ }^{\prime}, p_{\text {end }}{ }^{\prime}, j\right) \in S$ if and only if $\left.s \mu^{j}\right|_{p_{\text {start }}}=x$ and $p_{\text {end }}{ }^{\prime}$ is the minimal extension of $p_{\text {start }}{ }^{\prime}$.

First consider invariant (i). Clearly it is satisfied initially since $p_{\text {start }} p_{\text {end }}=$ $\varepsilon q=q=q p^{0}$. As induction hypothesis we obtain some minimal $k$ such that $\overline{p_{\text {start }} p_{\text {end }}}=q p^{k}$. After Steps (2) and (3) the invariant holds since neither $p_{\text {start }}$ nor $p_{\text {end }}$ are changed.

For Step (4) we first consider the case where Step (4-a) is not applied. Then we know that $\overline{p_{\text {end }}}=j p^{\prime}$ and $u=f\left(u_{1}, \ldots, u_{n}\right)$. Hence

$$
\begin{aligned}
p_{\text {start }} p_{\text {end }} & =\left(\overline{p_{\text {start }}} j\right) p^{\prime} \\
& =\overline{p_{\text {start }} p_{\text {end }}} \\
& =q p^{k}
\end{aligned}
$$

Minimality of $k$ is trivial, since $p_{\text {start }}$ is even longer than $\overline{p_{\text {start }}}$.
If in the execution of Step $\overline{(4)}$ case $(4-\mathrm{a})$ is applied we know that $\overline{p_{\text {end }}}=\varepsilon$ and hence, $\overline{p_{\text {start }}}=\overline{p_{\text {start }}} \varepsilon=\overline{p_{\text {start }} p_{\text {end }}}=q p^{k}$ by the IH. Moreover, $p_{\text {start }} p_{\text {end }}=$ $\overline{p_{\text {start }}} j p^{\prime}=\overline{p_{\text {start }}} p=q p^{k} p=q p^{k+1}$. Thus, $p_{\text {end }}$ is a minimal extension of $\underline{p_{\text {start }}}$ where minimality follows from the fact that $\overline{p_{\text {start }}}=\overline{p_{\text {start }}} j=q p^{k} j$ is not a prefix of $q p^{k}$.

Now for invariant (ii). Clearly it is satisfied initially. As induction hypothesis we have $\bar{u}=\left.s \mu^{\bar{i}}\right|_{\overline{p_{\text {star }}}}$.

After Steps (2) the invariant is satisfied, which can be proven directly by the IH since neither $u$, nor $i$, nor $p_{\text {start }}$ are changed.

After Step (3) we have

$$
\begin{aligned}
\underline{u} & =\bar{u} \mu \\
& =\left.s \mu^{\bar{i}}\right|_{\overline{p_{\text {start }}}} \mu \\
& =\left.s \mu^{\bar{i}+1}\right|_{\overline{p_{\text {start }}}} \\
& =\left.s \mu^{i}\right|_{p_{\text {start }}}
\end{aligned}
$$

using $\underline{p}_{\text {start }}=\overline{p_{\text {start }}}$.
For Step (4) we know that $\underline{p_{\text {start }}}=\overline{p_{\text {start }}} j, \bar{u}=f\left(u_{1}, \ldots, u_{n}\right)$, and $\underline{u}=u_{j}$. Hence

$$
\underline{u}=u_{j}=\left.\bar{u}\right|_{j}=\left.\left.s \mu^{\bar{i}}\right|_{\overline{p_{\text {statt }}}}\right|_{j}=\left.s \mu^{\bar{i}}\right|_{\overline{p_{\text {start }}} j}=\left.s \mu^{\underline{i}}\right|_{p_{\text {start }}}
$$

using the IH.
Invariant (iii) is satisfied initially, since $i=0$ and $S=\emptyset$. As induction hypothesis we have that $\left(x, p_{\text {start }}{ }^{\prime}, p_{\text {end }}{ }^{\prime}, j\right) \in \bar{S}$ iff $j<\bar{i},\left.s \mu^{j}\right|_{p_{\text {start }}}=x$, and $p_{\text {end }}{ }^{\prime}$ is the minimal extension of $p_{\text {start }}{ }^{\prime}$. Hence for Steps (2) and (4) that do neither change $i$ nor $S$, it trivially holds. Which leaves Step (3).

For Step (3) using the IH we only have to consider quadruples where the last component is $\bar{i}=\underline{i}-1$. Hence, we are interested in prefixes $p^{\prime}$ of $q p^{\omega}$ such that $\left.s \mu^{\bar{i}}\right|_{p^{\prime}}$ is a variable, together with their minimal extensions. Note that $\overline{p_{\text {end }}}$ is the minimal extension of $\overline{p_{\text {start }}}$ by invariant (i). From invariant (ii) and the condition of (3) we additionally know $\left.s \mu^{\bar{i}}\right|_{\overline{p_{\text {start }}}}=\bar{u}=x$. Hence, $\left(x, \overline{p_{\text {start }}}, \overline{p_{\text {end }}}, \bar{i}\right)$ has to be added to $\underline{S}$. Moreover, no other quadruples must be added to $\underline{S}$ since $\overline{p_{\text {start }}}$ is the unique prefix of $q p^{\omega}$ such that the subterm of $s \mu^{\bar{i}}$ at that position is a variable. Furthermore, the variable $x$ at that position is unique and also the minimal extension $\overline{p_{\text {end }}}$ of $\overline{p_{\text {start }}}$ is unique. This proves that invariant (iii) is establish after Step (3) which finishes the proof of the invariants.

With the help of the invariants it is now possible to prove soundness. So let Step (2) be applied. Hence, we have values $x, p_{\text {start }}, p_{\text {end }}, i$ where due to invariants (i) and (ii) there is some $k$ such that $p_{\text {start }} p_{\text {end }}=q p^{k}$ and $x=\left.s \mu^{i}\right|_{p_{\text {start }}}$. Moreover, since Step (2) is applied there also must be some element $\left(x, p_{\text {start }}^{\prime \prime}, p_{\text {end }}, i^{\prime \prime}\right) \in S$. Thus, by invariant (iii) we know that there is a $k^{\prime \prime}$ such that $p_{\text {start }}{ }^{\prime \prime} p_{\text {end }}=q p^{k^{\prime \prime}}$ and $x=\left.s \mu^{i^{\prime \prime}}\right|_{p_{\text {start }}}$. Since every time something is added to $S$, the value of the variable $i$ is increased we must have $i^{\prime \prime}<i$, so let $i=i^{\prime \prime}+j$ for some $j>0$.

Obviously, $p_{\text {start }}{ }^{\prime \prime} \leq p_{\text {start }}$ since the variable $p_{\text {start }}$ is only increased. But moreover, $p_{\text {start }}{ }^{\prime \prime}<p_{\text {start }}$ must be satisfied. Suppose $p_{\text {start }}{ }^{\prime \prime} \nless p_{\text {start }}$, hence $p_{\text {start }}=p_{\text {start }}{ }^{\prime \prime}$ and thus

$$
x=\left.s \mu^{i}\right|_{p_{\text {start }}}=\left.s \mu^{i^{\prime \prime}+j}\right|_{p_{\text {start }}}=\left.s \mu^{i^{\prime \prime}+j}\right|_{p_{\text {start }}^{\prime \prime}}=\left.s \mu^{i^{\prime \prime}}\right|_{p_{\text {start }}} \mu^{j}=x \mu^{j} .
$$

But then $x \mu^{j^{\prime}}$ will always be a variable regardless of $j^{\prime}$ and therefore, when traversing some term $s \mu^{j^{\prime}}$ along the path $q p^{\omega}$ one will reach a variable at position $p_{\text {start }}$ at the latest. This is in contradiction to the requirement $q p^{*} \subseteq$ $\bigcup_{j \in \mathbb{N}} \mathcal{P o s}\left(s \mu^{j}\right)$.

We continue to conclude from $p_{\text {start }}{ }^{\prime \prime}<p_{\text {start }}, p_{\text {start }} p_{\text {end }}=q p^{k}$, and $p_{\text {start }}{ }^{\prime \prime} p_{\text {end }}=$ $q p^{k^{\prime \prime}}$ that $k^{\prime \prime}<k$, so let $k=k^{\prime \prime}+k^{\prime}+1$. Since $p_{\text {end }}$ is the minimal extension of $p_{\text {start }}$ and since $k>0$ we know that $p_{\text {end }}$ must be a suffix of $p$, i.e., $p=p_{\text {pre }} p_{\text {end }}$ for some prefix $p_{\text {pre }}$ of $p$. Thus,

$$
p_{\text {start }} p_{\text {end }}=q p^{k^{\prime \prime}} p^{k^{\prime}} p=p_{\text {start }}{ }^{\prime \prime} p_{\text {end }} p^{k^{\prime}} p=p_{\text {start }}{ }^{\prime \prime} p_{\text {end }} p^{k^{\prime}} p_{\text {pre }} p_{\text {end }}
$$

and hence, by removing the final $p_{\text {end }}$ we obtain $p_{\text {start }}=p_{\text {start }}{ }^{\prime \prime} p_{\text {end }} p^{k^{\prime}} p_{\text {pre }}$. Thus,

$$
\left.x \mu^{j}\right|_{p_{\text {end }} p^{k^{\prime}}} p_{\text {pre }}=\left.\left.s \mu^{i^{\prime \prime}}\right|_{p_{\text {start }}^{\prime \prime}} \mu^{j}\right|_{p_{\text {end }} p^{k^{\prime}}} p_{\text {pre }}=\left.s \mu^{i}\right|_{p_{\text {start }} p_{\text {end }} p^{k^{\prime}} p_{\text {pre }}}=\left.s \mu^{i}\right|_{p_{\text {start }}}=x .
$$

Hence, each time one applies $\mu$ for $j$ more times on variable $x$ then at position $p_{\text {end }} p^{k^{\prime}} p_{\text {pre }}$ one sees $x$ again. Thus, each term $\left.s \mu^{i^{\prime \prime}} \mu^{j n}\right|_{p_{\text {start }}\left(p_{\text {end }} p^{k^{\prime}} p_{\text {pre }}\right)^{n}}$ is the variable $x$. Hence, to prove that $x$ infinitely often occurs in terms $s \mu^{j^{\prime}}$ along path $q p^{\omega}$ it suffices to show that all positions $p_{\text {start }}{ }^{\prime \prime}\left(p_{\text {end }} p^{k^{\prime}} p_{\text {pre }}\right)^{n}$ are prefixes of $q p^{\omega}$. To this end we show by a simple case-distinction on $n$ that $p_{\text {start }}{ }^{\prime \prime}\left(p_{\text {end }} p^{k^{\prime}} p_{\text {pre }}\right)^{n} p_{\text {end }} \in q p^{*}$. If $n=0$ then we just obtain $p_{\text {start }} p_{\text {end }}$ which is $q p^{k^{\prime \prime}} \in q p^{*}$. Otherwise,

$$
\begin{aligned}
p_{\text {start }}{ }^{\prime \prime}\left(p_{\text {end }} p^{k^{\prime}} p_{\text {pre }}\right)^{n} p_{\text {end }} & =p_{\text {start }}{ }^{\prime \prime} p_{\text {end }} p^{k^{\prime}} p_{\text {pre }}\left(p_{\text {end }} p^{k^{\prime}} p_{\text {pre }}\right)^{n-1} p_{\text {end }} \\
& =q p^{k^{\prime \prime}+k^{\prime}} p_{\text {pre }}\left(p_{\text {end }} p^{k^{\prime}} p_{\text {pre }}\right)^{n-1} p_{\text {end }} \\
& =q p^{k^{\prime \prime}+k^{\prime}}\left(p_{\text {pre }} p_{\text {end }} p^{k^{\prime}}\right)^{n-1} p_{\text {pre }} p_{\text {end }} \\
& =q p^{k^{\prime \prime}+k^{\prime}+(n-1)\left(k^{\prime}+1\right)+1} \in q p^{*} .
\end{aligned}
$$

Minimality of the returned values $i^{\prime}$ and $p_{\text {start }}{ }^{\prime}$ are then obvious by the com-
 algorithm idx in Fig. 4 to check whether $x$ occurs along $p^{\omega}$ in some term $t \mu^{j}$.

Theorem 4. The algorithm idx terminates and is sound.

Preconditions: $p \neq \varepsilon$ and $p^{*} \subseteq \bigcup_{i \in \mathbb{N}} \mathcal{P} \operatorname{os}\left(t \mu^{i}\right)$
Input: $x, t, p, \mu$
Output: $(j, q)$ if $j$ is minimal such that $\left.t \mu^{j}\right|_{q}=x$ where $q \in \mathcal{P o s}\left(t \mu^{j}\right)$ is prefix of $p^{\omega}$ or $\perp$, if there is no such $j$
(1) $j:=0, u:=t, p_{\text {start }}:=\varepsilon, p_{\text {end }}:=\varepsilon, S:=\varnothing$
(2) if $u=x$ then return $\left(j, p_{\text {start }}\right)$
(3) if $u=y \neq x$ and $\left(y, p_{\text {end }}\right) \in S$ then return $\perp$
(4) if $u=y \neq x$ and $\left(y, p_{\text {end }}\right) \notin S$ then $j:=j+1, S:=S \cup\left\{\left(y, p_{\text {end }}\right)\right\}, u:=u \mu$, goto (2)
(5) (a) if $p_{\text {end }}=\varepsilon$ then $p_{\text {end }}:=p$
(b) let $p_{\text {end }}=i p^{\prime}$ and $u=f\left(u_{1}, \ldots, u_{n}\right) ; u:=u_{i}, p_{\text {end }}:=p^{\prime}, p_{\text {start }}:=p_{\text {start }} i$, goto (2)

Fig. 4. $\mathrm{id} \times(x, t, p, \mu)$

Proof. Using the preconditions that $p^{*} \subseteq \bigcup_{i \in \mathbb{N}} \mathcal{P} \operatorname{os}\left(t \mu^{i}\right)$ and $p \neq \varepsilon$, termination of the algorithm in Fig. 3 can be shown as follows. First note that $u$ will always be a subterm of $t$ or of $y \mu$ for some $y \in \operatorname{Dom}(\mu)$. Moreover, $p_{\text {end }}$ will always be a suffix of $p$. As long as $u$ is not a variable, Step (5-b) makes sure that at some point a variable is reached (using the preconditions). Further note that Step (4) can only be used finitely often since $\mu$ has finite domain and since there are only finitely many possible values for $p_{\text {end }}$. Hence, eventually Step (2) or Step (3) must be applied and the algorithm stops.

To prove soundness we first prove the following invariants which are valid before Step (2).
(i) $p_{\text {start }} p_{\text {end }}=p^{k}$ for some $k$
(ii) $u=\left.t \mu^{j}\right|_{p_{\text {start }}}$
(iii) $\left(y,,_{-}\right) \in S$ iff $\left.t \mu^{i}\right|_{q_{i}}=y$ and $i<j$ where $q_{i}$ is the unique position s.t. $q_{i}<p^{\omega}$ and $q_{i} \in \mathcal{P o s}\left(t \mu^{i}\right)$.

The invariants are obviously satisfied after execution of Step (1). For Steps (2) and (3) there is nothing to prove. So, let us consider Step (4). Then invariant (i) directly follows from the IH since neither $p_{\text {start }}$ nor $p_{\text {end }}$ is changed. Moreover,

$$
\underline{u}=\bar{u} \mu=\left.t \mu^{\bar{j}}\right|_{\overline{p_{\text {start }}}} \mu=\left.t \mu^{\bar{j}} \mu\right|_{\overline{p_{\text {start }}}}=t \mu^{\underline{j}-\underline{p_{\text {start }}}}{ }
$$

shows that invariant (ii) is satisfied where we again used the IH. Finally, by using the IH of (iii) we only have to show that $\left.t \mu^{j-1}\right|_{q_{\underline{j}-1}}=\left.t \mu^{\bar{j}}\right|_{q_{\bar{j}}}=y$ where $q_{\bar{j}}^{-}$ is the maximal prefix of $p^{\omega}$ s.t. $q_{\bar{j}} \in \mathcal{P} \operatorname{os}\left(t \mu^{\bar{j}}\right)$. But this follows from IHs (i) and (ii) since $\overline{p_{\text {start }}}<p^{\omega}, \overline{p_{\text {start }}} \in \mathcal{P o s}\left(t \mu^{\bar{j}}\right)$, and $y=\bar{u}=\left.t \mu^{\bar{j}}\right|_{p_{\text {start }}}$. Due to the latter equality we see that $\overline{p_{\text {start }}}$ is maximal, and thus, $\overline{p_{\text {start }}}=q_{\bar{j}}$. Hence, we obtain $y=\left.t \mu^{\bar{j}}\right|_{\overline{p_{\text {start }}}}=\left.t \mu^{\bar{j}}\right|_{q_{\bar{j}}}$ and have therefore proven all invariants for Step (4).

It remains to prove the invariants for Step (5). Since neither $j$ nor $S$ are changed, there is nothing to show for invariant (iii). For invariant (i) there are two cases. If Step (5-a) is applied then the IH yields

$$
\underline{p_{\text {start }} p_{\text {end }}}=\left(\overline{p_{\text {start }}} i\right) p^{\prime}=\overline{p_{\text {start }}} \varepsilon\left(i p^{\prime}\right)=\overline{p_{\text {start }} p_{\text {end }}} p=p^{k} p=p^{k+1}
$$

where the intermediate value of $p_{\mathrm{end}}=p=i p^{\prime}$. Otherwise,

$$
\underline{p_{\text {start }} p_{\text {end }}}=\left(\overline{p_{\text {start }}} i\right) p^{\prime}=\overline{p_{\text {start }}}\left(i p^{\prime}\right)=\overline{p_{\text {start }} p_{\text {end }}}=p^{k}
$$

shows that invariant (i) is established. For proving invariant (iii) we know that $\underline{p_{\text {start }}}=\overline{p_{\text {start }}} i, \bar{u}=f\left(u_{1}, \ldots, u_{n}\right)$, and $\underline{u}=u_{i}$. Hence

$$
\underline{u}=u_{i}=\left.\bar{u}\right|_{i}=t \mu^{\bar{j}}\left|{\overline{p_{\text {start }}}}\right|_{i}=t \mu^{\bar{j}}\left|\frac{\bar{p}_{\text {start }} i}{}=t \mu^{\underline{j}}\right|_{\underline{p_{\text {start }}}}
$$

using the IH.
With the help of the invariants we now show that the function returns the correct result where we start with Step (2). By invariants (i) and (ii) we know that $x=u=\left.t \mu^{j}\right|_{p_{\text {start }}}$ where $p_{\text {start }}$ is a prefix of $p^{\omega}$. Thus, it remains to show minimality of $j$. So, suppose there is some $i<j$ such that $x=\left.t \mu^{i}\right|_{q_{i}}$ for some position $q_{i}<p^{\omega}$. Then by invariant (iii) we know that there is some entry $\left(x,{ }_{-}\right) \in S$. This contradicts the fact that the algorithm only inserts pairs into $S$ where the first component is not $x$.

We finally have to prove soundness of Step (3). So, suppose that the algorithm returns $\perp$ for some value of $j, p_{\text {end }}$, and $u=y$ whereas there are $j^{\prime}$ and $q_{j^{\prime}}<p^{\omega}$ such that $\left.t \mu^{j^{\prime}}\right|_{q_{j^{\prime}}}=x$. Note that due to invariant (iii) we know $j^{\prime} \geq j$ and by invariants (i) and (ii) we conclude $j^{\prime} \neq j$, and hence $j^{\prime}>j$.

Now consider a modification of the algorithm where Step (3) is missing and where Step (4) ignores the condition $\left(y, p_{\text {end }}\right) \notin S$. Then still all three invariants are satisfied. Hence, this modified algorithm can never terminate with $\perp$ and thus, must terminate using Step (2) immediately after the value of $j$ has been increased to $j^{\prime}$. We will show that this gives rise to a contradiction. The reason is that the values of $u=y$ and $p_{\text {end }}$ have already been encountered earlier during the run when inserting $\left(y, p_{\text {end }}\right)$ into $S$. And between this initial insertion and the application of Step (3) of the original algorithm, one has applied a sequence of steps which obviously does not include any application of Step (2). But since the upcoming steps of the modified algorithm are completely determined by $u$ and $p_{\text {end }}$, one will perform this sequence of steps over and over again. This contradicts the fact that the modified algorithm terminates by Step (2).

Note that both algorithms var $_{\infty}$ and idx become unsound if one only considers equal variables in the set $S$, but not equal positions $p_{\text {end }}$.

## 6 Empirical Results

In the following we first give some details about the actual implementation of our method and after that empirical results. To find loops, we use unfoldings as defined in [12], Section 3 (without any refinements mentioned in later sections). For efficiency reasons we restrict to non-variable positions. Further we do not use a combination of forward and backward unfoldings by default. Our basic method uses the following heuristic to decide the direction of the unfoldings:

For systems that are duplicating but whose inverse is non-duplicating we unfold backwards. For all other systems we unfold forwards.

In contrast to finding loops for full termination, for a specific strategy $\mathcal{S}$, we cannot always stop when a loop was found (since it could turn out to be no longer relevant when switching from full rewriting to $\mathcal{S}$-rewriting). Hence we compute a lazy list of potential loops that is checked one at a time corresponding to $\mathcal{S}$. If the checked loop is no $\mathcal{S}$-loop, the next loop is requested (and lazyness makes sure that the necessary computations are only done after the previous element dropped out).

For context-sensitive rewriting as well as outermost rewriting we reduce the search space by filtering the set of unfoldings after each iteration. In both cases we remove derivations containing rewrite steps that disobey the strategy.

As already mentioned in the introduction there is the transformational approach for both, context-sensitive rewriting [3] and outermost rewriting [13,14]. In both cases a problem for finding loops, is that the length of an existing loop may increase dramatically, or even worse, a loop is transformed into a nonlooping infinite derivation.

To evaluate our implementation in $\mathrm{T}_{\mathrm{T}} \mathrm{T}_{2}$ we used all 291 (214) outermost examples as well as all 109 (15) context-sensitive examples of version 5.0.2 of the Termination Problems Data Base. ${ }^{4}$ Here, in brackets the number of those TRSs is given, where outermost resp. context-sensitive termination has not already been proven. The results on these possibly non-terminating TRSs are as follows:

|  | outermost TRSs |  |  | CSRs |
| ---: | :---: | :---: | :---: | :---: |
|  | AProVE | TrafO | $\mathrm{T}^{\mathbf{T}} \mathrm{T}_{2}$ | $\mathrm{~T}^{\top} \mathrm{T}_{2}$ |
| NO score | 37 | 30 | 191 | 4 |
| avg. time (msec) | 6689 | 6772 | 340 | 38 |

For outermost rewriting we compare to the sum of non-termination proofs (NO score) achieved by AProVE and TrafO ${ }^{5}$ at the January 2009 termination competition. ${ }^{46}$ The success of our technique is clearly visible: $\mathrm{T}_{\mathrm{T}} \mathrm{T}_{2}$ was able to disprove termination of nearly $90 \%$ of all possible non-terminating TRSs, including all examples that could be handled by AProVE and TrafO (which use the transformational approaches of [14] and [13] respectively).

For context-sensitive rewriting we just give the NOs of $\mathrm{T}_{\top} T_{2}$ since we are not aware of any other tool that has disproven context-sensitive termination of a single TRS. Here, our implementation could at least solve one quarter of the potentially non-terminating TRSs.

The details of our experiments are available at http://cl-informatik. uibk.ac.at/~griff/experiments/lus.php.

[^4]
## 7 Conclusion and Future Work

To prove non-termination of rewriting under strategy $\mathcal{S}$, we first extended the notion of a loop to an $\mathcal{S}$-loop. An $\mathcal{S}$-loop is an $\mathcal{S}$-reduction with a strong regularity which admits the same infinite reduction as an ordinary loop does for full rewriting. Afterwards, we developed two novel procedures to decide whether a given loop is a context-sensitive loop or an outermost loop. It is easy to see that the conjunction of both procedures decides context-sensitive outermost loops.

Since [6] only describes a way to prove termination of Haskell programs, it might be an interesting future work to combine our technique for outermost loops with [6] to also disprove termination of Haskell programs.

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[^1]:    ${ }^{1}$ Another natural definition of an $\mathcal{S}$-loop would just require that $t(C, \mu)^{n} \rightarrow^{+}$ $t(C, \mu)^{n+1}$ are $\mathcal{S}$-derivations for all $n$. This alternative was already used in the setting of dependency pairs in [4, Footnote 6]. However, there are problems using this definition which are described in [15, Sect. 2].

[^2]:    $\overline{2}$ Note that in [15] one did not regard contexts, i.e., for an innermost loop one just required that all reductions $t_{i} \mu^{m} \rightarrow_{q_{i}} t_{i+1} \mu^{m}$ are innermost reductions. However, that definition of an innermost loop is equivalent to Def. 2 since innermost rewriting (denoted by $\xrightarrow{i}$ ) is monotonic. Thus, $t_{i} \mu^{m} \xrightarrow{i}_{q_{i}} t_{i+1} \mu^{m}$ iff $t_{i}(C, \mu)^{m} \stackrel{i}{p}_{p^{m}} q_{i} t_{i+1}(C, \mu)^{m}$.

[^3]:    ${ }^{3}$ Here we need the assumption $C \neq \square$. Otherwise, Rule (viii) would not terminate.

[^4]:    ${ }^{4}$ http://termcomp.uibk.ac.at
    ${ }^{5}$ http://www.win.tue.nl/~mraffels/trafo.html
    ${ }^{6}$ The numbers for $\mathrm{T} \mathrm{T}_{2}$ differ from those of the competition since the competition version did not feature the reduction of the search space which is described above.

