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# Generating invariants for non-linear loops by linear algebraic methods 

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#### Abstract

We present new computational methods that can automate the discovery and the strengthening of nonlinear interrelationships among the variables of programs containing non-linear loops, that is, that give rise to multivariate polynomial and fractional relationships. Our methods have complexities lower than the mathematical foundations of the previous approaches, which used Gröbner basis computations, quantifier eliminations or cylindrical algebraic decompositions. We show that the preconditions for discrete transitions can be viewed as morphisms over a vector space of degree bounded by polynomials. These morphisms can, thus, be suitably represented by matrices. We also introduce fractional and polynomial consecution, as more general forms for approximating consecution. The new relaxed consecution conditions are also encoded as morphisms represented by matrices. By so doing, we can reduce the non-linear loop invariant generation problem to the computation of eigenspaces of specific morphisms. Moreover, as one of the main results, we provide very general sufficient conditions allowing for the existence and computation of whole loop invariant ideals. As far as it is our knowledge, it is the first invariant generation methods that can handle multivariate fractional loops.


Keywords: Formal methods, Invariant generation, Linear algebra

## 1. Introduction

An invariant at a location of a program is an assertion true of any reachable program state associated to this location. We present a new method for non-linear invariant generation that addresses various deficiencies of other state-of-the-art methods. More generally, we provide mathematical techniques and design efficient algorithms to automate the discovery and strengthening of non-linear interrelationships among the variables of programs containing non-linear loops, which lead to multivariate polynomial and fractional relationships.

[^0]It is well-known that the automation and effectiveness of formal verification depend on the ease with which invariants can be automatically generated. Actually, the verification problem of safety properties, such as no null pointer deferenciation, buffer overflows, memory leak or outbounds, and array accesses, can be reduced to the problem of invariant generation [MP95]. Invariants are also essential to prove and establish liveness properties such as progress or termination [MP95]. Furthermore, the standard techniques [MP95] for program verification use invariant assertions to directly prove program properties, or to provide supporting lemmas that can be used to establish other safety and liveness properties. We look for invariants that strengthen what we wish to prove, and so allow us to establish the desired property. Also, they can provide precise over-approximations to the set of reachable states. Also, the weakest precondition method [Dij76, Flo67], the Floyd-Hoare inductive assertion technique [Flo67, Hoa69], and the standard ranking functions technique [MP95], all require loop invariants in order to establish correctness and so render the verification method completely automatic. Again, in order to establish termination verification, the standard ranking functions technique requires the automatic generation of invariants.

In order to generate loop invariants, one need to discover inductive assertions that hold at any step of the loop. An inductive assertion also holds at the first time the loop location is reached - this is the initiation condition - and it is also preserved under instructions that cycle back to the loop location - this being the consecution condition. If we choose transition systems as the representation model and automata as the computational model, we can say that the invariant holds in the initial state of the system - the initial condition - and that every possible transition preserves it-the consecution conditions. In other words, the invariant holds in any possible reachable state.

In the case of loops describing a linear system, Farka's lemma [Sch86] can be used to encode the conditions for linear invariants. On the other hand, for non-linear invariants, the difficulty of automatic generation remains very challenging. By today known methods, they require a large number of Gröbner bases computations [SSM04b], first-order quantifier eliminations [Wei97, Col75], or cylindrical algebraic decompositions [CXYZ07]. Invariants can also be computed as fixed points on ideals, using fixed point techniques [RCK07a], abstract interpretation frameworks [CC92, CC77], and Gröbner bases constructions. Abstract interpretation introduces imprecision, and widening operators must be provided manually by the user in order to assure termination. A too coarse abstraction would limit these approaches to trivial invariants in the presence of non-linear loops. Other methods [KJ06, Kov08] attempt to generate invariants from a restricted class of P-solvable loops. These methods use techniques from algebra and combinatorics, like Gröbner bases [JKP06], variable elimination, algebraic dependencies and symbolic summation, and so also incur in high computational complexities.

More recent approaches have been constraint-based [SSM04b, RCK07a, Kap04, RCK07b, SSM04a, GT08, PJ04]. In these cases, a candidate invariant with a fixed degree and unknown parametric coefficients, i.e., a template form, is proposed as the target invariant to be generated. The conditions for invariance are then encoded, resulting in constraints on the unknown coefficients whose solutions yield invariants. One of the main advantage of such constraint-based approaches is that they are goal-oriented. The main challenge for these techniques remains in the fact that they still require a high number of Gröbner bases [Buc96] computations, first-order quantifier elimination [Wei97, Co175], cylindrical algebraic decomposition [CXYZ07], or abstraction operators. And known algorithms for those problems are, at least, of double exponential complexity.

Despite tremendous progress over the years [SSM04b, BBGL00, RCK07a, BLS96, CXYZ07, Kov08, KJ06, Cou05, MOS02, RCK07b, GT08, Tiw08, PC08], the problem of loop invariant generation remains very challenging for non-linear discrete systems. In this work we present new methods for the automatic generation of loop invariants for non-linear systems. As will be seen, these methods give rise to more efficient algorithms, with much lower complexity in space and time. We develop the new methods by first extending our previous work on non-linear non-trivial invariant generation for discrete programs with nested loops and conditional statements, [RMM08b, RMM10].

We can summarize our contributions as follows:

- We do not need to start with candidate invariants that generate intractable solving problems. Instead, we show that the preconditions for discrete transitions can be viewed as morphisms over a vector space of degree bounded by polynomials which can, thus, be suitably represented by matrices.
- We introduce more general forms for approximating consecution, called fraction and polynomial consecution. The new relaxed consecution requirements are also encoded as morphisms, represented by matrices with terms that are the unknown coefficients used to approximate the consecution conditions. As far as it is our knowledge, these are the first methods that can effectively handle multivariate fractional systems.
- We succeed in reducing the non-linear loop invariant generation problem to the computation of eigenspaces or nullspaces of specific endomorphisms. We provide general sufficient conditions guaranteeing the existence and allowing the computation of invariant ideals. The unknown coefficients appearing in the matrices used to approximate the consecution conditions are assigned in order to insure that the nullspaces generated are not trivial ones. Taking into consideration the specific type of matrices we are manipulating, we determine for which values of the coefficients their ranks are minimal. Our decision procedure for those assignments is very simple and efficient. At each step of the assignments, we echelon the matrices by making the highest term of one column to vanish.
- Our approach does not generate an invariant at a time. Instead we generate an ideal of invariants-an infinite structure-by computing the basis of a specific vector space giving rise to provable, inductive invariants. This could also be used by existing approaches dealing with the generations of such vector spaces of inductive invariants [Cou05, RCK07b, Kov08].
- Our technique comprises three computational steps, each of polynomial time complexity. In contrast, the most recent and best performing constraint-based approaches can be summarized in three main steps, with each of these steps inducing a number of computations that are of double exponential time complexity. Further, as soon as the loop contains non-linear instructions, the constraints considered at the final step gives rise to systems of non-linear equations, rendering unfeasible their resolution; see Sect. 4.3. We, by contrast, propose a computational method of much lower time complexity than other present approaches based on fixed point computation, or on constraint-based approaches.
- We present some preliminary experimental results. For that, we used Sage [SJ05] with interfaces written in Python, in order to be able to access other mathematical packages.
- We incorporate a strategy that attains optimal degree bounds for candidate invariants. We also note that our existence results and methods can be reused in other approaches in order to reduce their time complexity, since they can reduce the number of Gröbner basis computations or quantifier eliminations, for example.
Example 1 (Motivational Example). Consider the following program loop:

```
while (...){
    x := x*y + x;
    ...
    y := y^2;
}
```

Present constraint-based static program analysis techniques are facing some difficulties in producing any conclusion that could be somehow related to the values of the variables $x$ and $y$, given that the semantic of the two instructions inside the loop relies on non-linear arithmetic. Such non-linearities are presently recognized by industry and academia as a critical bottleneck for automatic program verification and static program analysis.

In present standard approaches for invariant generation for non-linear loops, the loop instructions are first used in order to form varieties, to build associated algebraic assertions and an ideal $I$. Then, they compute a Gröbner basis $G$ for $I$. Next, they postulate a template polynomial $Q$, i.e., a polynomial with unknown coefficients, as a candidate invariant, and proceed by performing a reduction of $Q$ by $G$ in order to obtain its reduced normal form $N F_{G}(Q)$. An important obstacle faced at this point is that all known algorithms for computing Gröbner basis and for constructing the normal form reduction $N F_{G}(Q)$ are of doubly exponential time complexity. Having the normal form $N F_{G}(Q)$, they generate the set of candidate invariant constraints in the form of a system of equations by letting $N F_{G}(Q)=0$, and then they attempt to solve it directly. But we show that as soon as the loop contains a non-linear instruction, the constraints obtained in their final step lead to systems of non-linear equations in unknown parameters, which remains untractable in practice (see Sect. 4.3). For more details on the limitations of such techniques, illustrated on this same motivational example, see Example (6), in Sect. 4.3.

In this article, we introduce new symbolic techniques with fast numerical approaches that can be used in these situations. Our techniques have fewer computational steps. We first compute some specific matrix $M$ obtained directly from the loop instructions. We then generate a matrix $L$ that we use to approximate the consecution condition. Matrix $L$ contains some fixed parameters so as to guarantee that such nullspaces are not empty. We note that the unknown coefficients do not play the role of templates. Rather, they are introduced to allow us to reduce the rank of the matrix $M-L$, thus leading to a non-trivial nullspace. Based on our theoretical contributions, we know that the nullspace of $M-L$ provides us with a non-trivial vector space of inductive invariants. In the example at hand, our method directly computes $\left\{x^{2}, x * y-x, y^{2}-2 y+1\right\}$ as a basis for a vector space of invariants, and all elements in this space provide non-trivial invariants. We thus obtain an ideal for non-trivial inductive invariants. In other words, for all $G_{1}, G_{2}, G_{3} \in \mathbb{R}[x, y]$, we would get $G_{1}(x, y)\left(x^{2}\right)+G_{2}(x, y)(x y-x)+G_{3}(x, y)\left(y^{2}-2 y+1\right)=0$ as an inductive invariant. Take, for instance, the initial step $\left(y=y_{0}, x=1\right)$. A possible invariant is, then, $y_{0}\left(1-y_{0}\right) x^{2}+x y-x+y^{2}-2 y+1=0$. By taking two elements of this basis, one could generates inductive invariants holding for any type of initial conditions on the variables. Such invariants are beyond the reach of other current invariant generation techniques. In Sect. 5.4, we make explicit all the computational steps of our method.

In Sect. 2 we present ideals of polynomials and their possible interactions with inductive assertions. In Sect. 3 we introduce new consecution conditions, and extend them to fractional systems. In Sect. 4 we consider linear loops, and present results for the existence of non-trivial invariants in these settings. We also recast the problem in term of linear algebra, and present a complete decision procedure for the automatic generation of non-trivial non-linear invariants. In Sect. 5 we extend our method to non-linear loops. In Sect. 6 we propose a strategy to obtain optimal degree bounds. In Sect. 7 we provide a complete generalization by considering loops describing multivariate fractional systems, and in Sect. 8 we show how to handle conditions and nested loops. Section 9 exposes some preliminary experimental results, and Sect. 10 contains a discussion. We conclude in Sect. 11. The Appendix contains proofs for all theorems, lemmas and corollaries stated in this article. Further examples can be found in companion technical reports and other articles [RMM08a, RMM08b, RMM10, RM11a, RM11b].

## 2. Polynomial ideals and inductive assertions

We will use the following notations. Let $\mathbb{K}$ be a field. The ring of multivariate polynomials over the set of variables $\left\{X_{1}, . ., X_{n}\right\}$ with coefficients in $\mathbb{K}$ will be indicated by $\mathbb{K}\left[X_{1}, . ., X_{n}\right]$. We will denote by $\mathbb{R}_{d}\left[X_{1}, . ., X_{n}\right]$ the vector space of multivariate polynomials of degree at most $d$ over the set of real variables $\left\{X_{1}, \ldots, X_{n}\right\}$. We will write $\operatorname{Vect}\left(v_{1}, \ldots, v_{n}\right)$ for the vector space generated by a basis $\left(v_{1}, \ldots, v_{n}\right)$. The dimension of a subspace $W \subseteq \operatorname{Vect}\left(v_{1}, \ldots, v_{n}\right)$ is written $\operatorname{Dim}(W)$. Clearly, $\operatorname{Dim}\left(\operatorname{Vect}\left(v_{1}, \ldots, v_{n}\right)\right)=n$. The vector space of all matrices over a filed $\mathbb{K}$ will be denoted by $\mathcal{M}(m, n, \mathbb{K})$. Let $M \in \mathcal{M}(m, n, \mathbb{K})$ be the matrix representation of a morphism over a vector space. Its kernel, or nullspace, is the set $\operatorname{Ker}(M)=\left\{v \in \mathbb{K}^{n} \quad \mid \quad M \cdot v=0_{\mathbb{K}^{m}}\right\}$. The kernel of $M$ is said to be trivial if it contains only the zero vector. The rank of $M$, denoted $\operatorname{Rank}(M)$, is the dimension of the subspace $\left\{M \cdot v \in \mathbb{K}^{m} \mid v \in \mathbb{K}^{n}\right\}$. Alternatively, it is the number of linearly independent columns or rows of the matrix. We know that $\operatorname{Rank}(M)+\operatorname{Dim}(\operatorname{Ker}(M))=n$. An eigenvalue of $M$ is a scalar $\lambda \in \mathbb{K}$ such that $M \cdot v=\lambda v$ for some nonzero vector $v$. The set $\left\{v \in \mathbb{K}^{n} \mid M \cdot v=\lambda v\right\}$ is the eigenspace associated to an eigenvalue $\lambda$. To compute the basis of eigenspaces and nullspaces, we use well-known state-of-the-art algorithms, such as those that Sage or Mathematica provide. To solve equations of degree less than 5 one could consider the classical Lagrange resolvents [Lan02, AV97] method. A primed $x^{\prime}$ will refer to the next state value of a variable $x$, after a transition is taken. If $V$ is a set of variables, then $V^{\prime}$ is the set of all primed variables in $V$.

### 2.1. Polynomial ideals

Definition 1 An ideal is any set $I \subseteq \mathbb{K}\left[X_{1}, . ., X_{n}\right]$ such that

- It is closed under addition. In other words, if $P, Q \in I$ then $P+Q \in I$;
- It is closed under multiplication by any element in $\mathbb{K}\left[X_{1}, \ldots, X_{n}\right]$, i.e., if $P \in I$ and $Q \in K\left[X_{1}, . ., X_{n}\right]$ then $P Q \in I$;
- It includes the null polynomial, i.e. $0_{\mathbb{K}\left[X_{1}, \ldots, X_{n}\right]} \in I$.

Let $E \subseteq \mathbb{K}\left[X_{1}, . ., X_{n}\right]$ be a set of polynomials. The ideal generated by $E$ is the set of finite sums

$$
(E)=\left\{\sum_{i=1}^{k} P_{i} Q_{i} \mid P_{i} \in \mathbb{K}\left[X_{1}, \ldots, X_{n}\right], \quad Q_{i} \in E, k \geq 1\right\}
$$

Definition 2 A set of polynomials $E$ is said to be a basis of an ideal $I$ if $I=(E)$.
By the Hilbert Basis Theorem, we know that all ideals have a finite basis. Let $I$ be an ideal and $Q$ a polynomial. The important question of knowing if $Q$ belongs to $I$ is known as the Ideal Membership Problem. In order to decide membership, one first computes the normal form of $Q$ by performing polynomial reductions according to $I$. If the resulting normal form is the null polynomial we can conclude that $Q \in I$. A Gröbner basis of $I$ guarantees the confluence and termination of those polynomial reductions.

### 2.2. Inductive assertions and invariants

The contribution and novelty in our approach clearly set it apart from [SSM04b] as their constraint-based techniques require several Gröbner basis computations and also require solving non-linear problems for each location. Nevertheless, they introduce a useful formalism to treat programs loops, and we start from similar definitions for transitions systems, inductive invariants and consecution conditions. We will use transition systems as representation of imperative programs and automata as their computational models.

Definition 3 A transition system is given by $\left\langle V, L, \mathcal{T}, l_{0}, \Theta\right\rangle$, where

- $V$ is a set of variables,
- $L$ is a set of locations and $l_{0} \in L$ is the initial location.
- A transition $\tau \in \mathcal{T}$ is given by a tuple $\left\langle l_{\text {pre }}, l_{\text {post }}, \rho_{\tau}\right\rangle$, where $l_{\text {pre }}$ and $l_{\text {post }}$ name the pre- and post- locations of $\tau$, and the transition relation $\rho_{\tau}$ is a first-order assertion over $V \cup V^{\prime}$.
- $\Theta$ is the initial condition, given as a first-order assertion over $V$.

The transition system is affine when $\rho_{\tau}$ is an affine form, and it is algebraic when $\rho_{\tau}$ is an algebraic form.
Definition 4 Let $W$ be a transition system. An invariant at location $l \in L$ is an assertion over $V$ which holds at all states reaching location $l$. An invariant of $W$ is an assertion over $V$ that holds at all locations.

Given our representational and computational models, we want to say that an invariant holds in the initial state of the system, a condition that will be guaranteed by an initial condition. We also want to say that every possible transition preserves the invariant, when specific consecution conditions hold. In order to generate loop invariants one needs to discover inductive assertions.
Definition 5 Let $W=\left\langle V, L, \mathcal{T}, l_{0}, \Theta\right\rangle$ be a transition system and let $\mathbb{D}$ be an assertion domain. An assertion map for $W$ is a map $\eta: L \rightarrow \mathbb{D}$. We say that $\eta$ is inductive if and only if the following conditions hold:

- Initiation: $\Theta \models \eta\left(l_{0}\right)$
- Consecution: For all $\tau$ in $\mathcal{T}$ s.t. $\tau=\left\langle l_{i}, l_{j}, \rho_{\tau}\right\rangle$ we have $\eta\left(l_{i}\right) \wedge \rho_{\tau} \models \eta\left(l_{j}\right)^{\prime}$.

We know that if $\eta$ is an inductive assertion map then $\eta(l)$ is an invariant at $l$ for $W$ [Flo67].

## 3. New continuous consecution conditions

In this section we treat discrete transitions by extending and adapting our previous work on loop invariant generation for discrete programs [RMM08a, RMM08b, RMM10]. We also consider discrete transitions that are part of connected components and circuits, thus generalizing the case of simple propagations.

First, we show how to encode continuous consecution conditions.
Definition 6 Consider a transition system $W=\left\langle V, L, \mathcal{T}, l_{0}, \Theta\right\rangle$. Let $\tau=\left\langle l_{i}, l_{j}, \rho_{\tau}\right\rangle$ be a transition in $\mathcal{T}$ and let $\eta$ be an algebraic inductive map with $\eta\left(l_{i}\right) \equiv\left(P_{\eta}\left(X_{1}, . ., X_{n}\right)=0\right)$ and $\eta\left(l_{j}\right) \equiv\left(P_{\eta}^{\prime}\left(X_{1}, . ., X_{n}\right)=0\right)$ where $P_{\eta}$ is a multivariate polynomial in $\mathbb{R}\left[X_{1}, . ., X_{n}\right]$ such that it has null values at $l_{i}$ and at $l_{j}$, i.e., before and after taking the
transition. Note that this does not imply that $P_{\eta}$ is the null polynomial. We identify the following notions when encoding continuous consecution conditions:

- We say that $\eta$ satisfies a Fractional-scale consecution for $\tau$ if and only if there exists a multivariate fractional $\frac{T}{Q}$ such that $\rho_{\tau} \vDash\left(P_{\eta}\left(X_{1}^{\prime}, \ldots, X_{n}^{\prime}\right)-\frac{T}{Q} P_{\eta}\left(X_{1}, \ldots, X_{n}\right)=0\right)$. We also say that $P_{\eta}$ is a $\frac{T}{Q}$-scale discrete invariant.
- We say that $\eta$ satisfies a Polynomial-scale consecution for $\tau$ if and only if there exists a multivariate polynomial $T$ such that $\rho_{\tau} \vDash\left(P_{\eta}\left(X_{1}^{\prime}, \ldots, X_{n}^{\prime}\right)-T P_{\eta}\left(X_{1}, . ., X_{n}\right)=0\right)$. We also say that $P_{\eta}$ is a polynomial-scale and a $T$-scale discrete invariant.
- We say that $\eta$ satisfies a Constant-scale consecution for $\tau$ if and only if there exists a constant $\lambda$ such that $\rho_{\tau} \vDash\left(P_{\eta}\left(X_{1}^{\prime}, \ldots, X_{n}^{\prime}\right)-\lambda P_{\eta}\left(X_{1}, \ldots, X_{n}\right)=0\right)$. We also say that $P_{\eta}$ is a constant-scale, or a $\lambda$-scale discrete invariant.

Constant-scale consecution encodes the fact that the numerical value of the polynomial $P_{\eta}$, associated with assertion $\eta\left(l_{i}\right)$, is given by $\lambda$ times its numerical value throughout the transition $\tau$. Polynomial-scale consecution encodes the fact that the numerical value of the polynomial $P_{\eta}$, associated with assertion $\eta\left(l_{i}\right)$, is given by $T$ times its numerical value throughout the transition $\tau$, where $T$ is a polynomial in $\mathbb{R}\left[X_{1}, \ldots, X_{n}\right]$. Such $T$ polynomials can be understood as template multiplicative factors, that is, they are polynomials with unknown coefficients. We are able to handle the general case when the loop describes a multivariate fractional system with Fractionalscale consecution. Fractional-scale consecution encodes the fact that the numerical value of the polynomial $P_{\eta}$, associated with assertion $\eta\left(l_{i}\right)$, is given by $\frac{T}{Q}$ times its numerical value throughout the transition $\tau$. The fractional $\frac{T}{Q}$ can contain unknown coefficients. As can be seen, the consecution conditions are relaxed when going from constant to fractional scaling.

## 4. Discrete transition and affine systems

In this section we treat constant-scale consecution encodings. Consider a transition systems corresponding to the loop $\tau=\left\langle l_{i}, l_{i}, \rho_{\tau}\right\rangle$ and its affine transition relation

$$
\rho_{\tau} \equiv\left[\begin{array}{c}
X_{1}^{\prime}=L_{1}\left(X_{1}, \ldots, X_{n}\right)  \tag{1}\\
\vdots \\
X_{n}^{\prime}=L_{n}\left(X_{1}, \ldots, x_{n}\right)
\end{array}\right] .
$$

Here, the loop instructions are affine or linear forms $L_{i}\left(X_{1}, \ldots, X_{n}\right)=\sum_{k=1}^{n} c_{i, k-1} X_{k}+c_{i, k}, 1 \leq i \leq n$.

### 4.1. Generating $\lambda$-scale invariants

We have the following $\lambda$-scale invariant characterization.
Theorem 1 Consider a transition system corresponding to a loop $\tau$ as described in Eq. (1). A polynomial $Q$ in $\mathbb{R}\left[X_{1}, \ldots, X_{n}\right]$ is a $\lambda$-scale invariant for constant-scale consecution with parametric constant $\lambda \in \mathbb{R}$ for $\tau$ if and only if $Q\left(L_{1}\left(X_{1}, . ., X_{n}\right), . ., L_{n}\left(X_{1}, . ., X_{n}\right)\right)=\lambda Q\left(X_{1}, . ., X_{n}\right)$.

Let the degree of $Q \in \mathbb{R}\left[X_{1}, . ., X_{n}\right]$ be $r$. We show that for good choices of $\lambda$ there always exists such a $\lambda$-invariant that is also non-trivial. We note that $Q\left(L_{1}\left(X_{1}, \ldots, X_{n}\right), \ldots, L_{n}\left(X_{1}, \ldots, X_{n}\right)\right)$ is also of degree $r$ because all $L_{i}$ 's are of degree 1 . Recasting the situation and $\rho_{\tau}$ into linear algebra, consider the morphism

$$
\mathscr{M}:\left\{\begin{aligned}
\mathbb{R}_{r}\left[X_{1}, \ldots, X_{n}\right] & \rightarrow \mathbb{R}_{r}\left[X_{1}, \ldots, X_{n}\right] \\
Q\left(X_{1}, \ldots, X_{n}\right) & \mapsto Q\left(L_{1}\left(X_{1}, \ldots, X_{n}\right), \ldots, L_{n}\left(X_{1}, . ., X_{n}\right)\right) .
\end{aligned}\right.
$$

This is indeed an endomorphism because all $L_{i}$ 's are of degree 1 . Let $M$ be its matrix representation in the canonical basis of $\mathbb{R}_{r}\left[X_{1}, ., X_{n}\right]$. First, we show how we can build matrix $M$.

Example 2 Consider the following loop $\rho_{\tau}=\left[\begin{array}{c}x_{1}^{\prime}=2 x_{1}+x_{2}+1 \\ x_{2}^{\prime}=3 x_{2}+4\end{array}\right]$. We have two polynomials of degree 1 , in two variables. They are $L_{1}\left(x_{1}, x_{2}\right)=2 x_{1}+x_{2}+1$, and $L_{2}\left(x_{1}, x_{2}\right)=3 x_{2}+4$. Consider the associated endomorphism $\mathscr{M}$ from $\mathbb{R}_{2}\left[x_{1}, x_{2}\right]$ to $\mathbb{R}_{2}\left[x_{1}, x_{2}\right]$. We want to obtain an associated matrix $M$ for it. For that, we can use $B_{1}=$ $\left(x_{1}^{2}, x_{1} x_{2}, x_{2}^{2}, x_{1}, x_{2}, 1\right)$ as a basis for $\mathbb{R}_{2}\left[x_{1}, x_{2}\right]$ and compute $\mathscr{M}(P)$ for all elements $P$ in $B_{1}$, expressing the results in the same basis. For the first column of $M$ we consider $P\left(x_{1}, x_{2}\right)=x_{1}^{2}$ as the first element of $B_{1}$, and compute $\mathscr{M}(P)=P\left(L_{1}\left(x_{1}, x_{2}\right), L_{2}\left(x_{1}, x_{2}\right)\right)$, which is expressed in $B_{1}$ as

$$
M=\left(\begin{array}{llllll}
4 & M\left(x_{1}^{2}\right)=4 x_{1}^{2}+4 x_{1} x_{2}+1 x_{2}^{2}+4 x_{1}+2 x_{2}+1 \times 1 \\
4 & 6 & 0 & 0 & 0 & 0 \\
1 & 3 & 9 & 0 & 0 & 0 \\
4 & 8 & 0 & 2 & 0 & 0 \\
2 & 7 & 24 & 1 & 3 & 0 \\
1 & 4 & 16 & 1 & 4 & 1
\end{array}\right)
$$

This concludes the example.
Now, let $Q \in \mathbb{R}\left[X_{1}, \ldots, X_{n}\right]$ be a $\lambda$-scale invariant for constant-scale consecution with parametric constant $\lambda \in \mathbb{R}$ for a given system defined by $L_{1}, . ., L_{n} \in \mathbb{R}\left[X_{1}, . ., X_{n}\right]$. By Theorem 1 , we have

$$
Q\left(L_{1}\left(X_{1}, \ldots, X_{n}\right), . ., L_{n}\left(X_{1}, \ldots, X_{n}\right)\right)=\lambda Q\left(X_{1}, \ldots, X_{n}\right)
$$

Using the associated endomorphism $\mathscr{M}$, we have:

$$
\begin{aligned}
Q\left(L_{1}\left(X_{1}, \ldots, X_{n}\right), \ldots, L_{n}\left(X_{1}, \ldots, X_{n}\right)\right) & =\lambda Q\left(X_{1}, \ldots, X_{n}\right) & & \Leftrightarrow \\
\mathscr{M}(Q) & =\lambda Q & & \Leftrightarrow \\
\mathscr{M}(Q) & =\lambda \mathscr{I}(Q) & & \Leftrightarrow \\
(\mathscr{M}-\lambda \mathscr{I})(Q) & =0_{\mathbb{R}\left[X_{1}, \ldots, X_{n}\right]} & & \Leftrightarrow \\
Q & \in \operatorname{Ker}(M-\lambda I), & & r \text { 俍 }
\end{aligned}
$$

where $\mathscr{I}$ is the identity endomorphism and $I$ is the associated identity matrix in $\mathbb{R}_{r}\left[X_{1}, . . X_{n}\right]$. Hence, $\lambda$ must be an eigenvalue of $M$ if we want to find a non null $\lambda$-invariant whose coefficients will be those of an eigenvector.

We can now state the following theorem.
Theorem 2 A polynomial $Q$ of $\mathbb{R}_{r}\left[X_{1}, \ldots, X_{n}\right]$ is $\lambda$-invariant for constant-scale consecution if and only if there exists an eigenvalue $\lambda$ of $M$ such that $Q$ belongs to the eigenspace corresponding to $\lambda$.

We also notice that, by construction, the last column of $M$ is always $(0, \ldots, 0,1)^{\top}$. Thus 1 is always an eigenvalue of $M$ with a corresponding eigenvector which leads to the trivial $\lambda$-invariant $Q\left(X_{1}, \ldots, X_{n}\right)=a$, where $a$ is the coefficient of the constant term. Eigenvalue 1 always gives the constant polynomial as a $\lambda$-invariant, but it might give better invariants for other eigenvectors if $\operatorname{dim}(\operatorname{Ker}(M-\lambda I)) \geq 2$, as we will see in the sequel.

Example 3 Looking at the eigenvalues of the matrix $M$ in Example 2, if we fix $\lambda$ to be 4 we get that the corresponding eigenspace is generated by the vector $(1,-2,1,-6,6,9)^{\top}$. Interpreted in the canonical basis of $\mathbb{R}\left[x_{1}, x_{2}\right]$, the associated 4-invariant is $Q\left(x_{1}, x_{2}\right)=1 x_{1}^{2}-2 x_{1} x_{2}+x_{2}^{2}-6 x_{1}+6 x_{2}+9$.

We first treat the general case where the transition system has only two variables. We will look for a $\lambda$-invariant $Q$ of degree two. Let

$$
\rho_{\tau}=\left[\begin{array}{l}
x_{1}^{\prime}=c_{1,0} x_{1}+c_{1,1} x_{2}+c_{1,2} \\
x_{2}^{\prime}=c_{2,0} x_{1}+c_{2,1} x_{2}+c_{2,2}
\end{array}\right] .
$$

Recall that we must solve the equation $Q\left(c_{1,0} X_{1}+c_{1,1} X_{2}+c_{1,2}, c_{2,0} X_{1}+c_{2,1} X_{2}+c_{2,2}\right)=\lambda Q\left(X_{1}, X_{2}\right)$. Thus, for $M$ we get the following matrix:

$$
\left(\begin{array}{cccccc}
c_{1,0} & c_{1,0} c_{2,0} & c_{2,0} & 0 & 0 & 0 \\
2 c_{1,0} c_{1,1} & c_{1,0} c_{2,1}+c_{1,1} c_{2,0} & 2 c_{2,0} c_{2,1} & 0 & 0 & 0 \\
c_{1,1}^{2} & c_{1,1} c_{2,1} & c_{2,1}^{2} & 0 & 0 & 0 \\
2 c_{1,0} c_{1,2} & c_{1,0} c_{2,2}+c_{1,2} c_{2,0} & 2 c_{2,0} c_{2,2} & c_{1,0} & c_{2,0} & 0 \\
2 c_{1,1} c_{1,2} & c_{1,1} c_{2,2}+c_{1,2} c_{2,1} & 2 c_{2,1} c_{2,2} & c_{1,1} & c_{2,1} & 0 \\
c_{1,2}^{2} & c_{1,2} c_{2,2} & c_{2,2}^{2} & c_{1,2} & c_{2,2} & 1
\end{array}\right)
$$

We see that the last column is as predicted, plus the matrix is block diagonal. Thus its characteristic polynomial is $P(\lambda)=(1-\lambda) P_{1}(\lambda) P_{2}(\lambda)$, with $P_{1}$ being the characteristic polynomial of

$$
\left(\begin{array}{ll}
c_{1,0} & c_{2,0} \\
c_{1,1} & c_{2,1}
\end{array}\right),
$$

and $P_{2}$ being the characteristic polynomial of

$$
\left(\begin{array}{ccc}
c_{1,0}{ }^{2} & c_{1,0} c_{2,0} & c_{2,0}{ }^{2} \\
2 c_{1,0} c_{1,1} & c_{1,0} c_{2,1}+c_{1,1} c_{2,0} & 2 c_{2,0} c_{2,1} \\
c_{1,1}^{2} & c_{1,1} c_{2,1} & c_{2,1}^{2}
\end{array}\right) .
$$

Here $P_{2}$ is of degree 3 and has at least one real root. This root can be computed by the Lagrange resolvent method. Choosing $\lambda$ to be this root, the corresponding eigenvectors will give non-trivial $\lambda$-invariants of degree two, since at least one of the coefficients of the monomials $x_{1}^{2}, x_{1} x_{2}$ and $x_{2}^{2}$ must be non null for such an eigenvector.
Corollary 1 Let $M$ be the matrix introduced in this section. The problem of finding a non-trivial $\lambda$-invariant is decidable if one of the following assertions is true:

- $M$ is block triangular (with $4 \times 4$ blocks or less),
- The eigenspace associated with eigenvalue 1 is of dimension greater than 1.


### 4.2. Intersection with initial hyperplanes

Let $Q \in \mathbb{R}_{r}\left[X_{1}, \ldots, X_{n}\right]$ be a $\lambda$-invariant for constant-scale consecution, that is,

$$
Q\left(L_{1}\left(X_{1}, . ., X_{n}\right), . ., L_{n}\left(X_{1}, . ., X_{n}\right)\right)=\lambda Q\left(X_{1}, . ., X_{n}\right) .
$$

Now let $u_{1}, \ldots, u_{n}$ be the initial values of $X_{1}, \ldots, X_{n}$. For the initial step we need $Q\left(u_{1}, \ldots, u_{n}\right)=0$. We have $P \mapsto P\left(u_{1}, \ldots, u_{n}\right)$ as a linear form in $\mathbb{R}_{r}\left[X_{1}, \ldots, X_{n}\right]$. Hence initial values correspond to a hyperplane in $\mathbb{R}_{r}\left[X_{1}, \ldots, X_{n}\right]$, given by the kernel of $P \mapsto P\left(u_{1}, \ldots, u_{n}\right)$. If we add the initiation step, $Q\left(X_{1}, \ldots, X_{n}\right)=0$ will be an inductive invariant (see Definition 4) if and only if there exists an eigenvalue $\lambda$ of $M$ such that $Q$ belongs to the intersection of the eigenspace corresponding to $\lambda$ and the hyperplane $Q\left(u_{1}, \ldots, u_{n}\right)=0$.
Theorem 3 A polynomial $Q$ in $\mathbb{R}_{r}\left[X_{1}, \ldots, X_{n}\right]$ is an inductive invariant for the affine loop (see Definition 5) with initial values ( $u_{1}, \ldots, u_{n}$ ) if and only if there is an eigenvalue $\lambda$ of $M$ such that $Q$ is in the intersection of the eigenspace of $\lambda$ and the hyperplane $Q\left(u_{1}, \ldots, u_{n}\right)=0$.

In the following corollary, we state an important result.
Corollary 2 There will be a non-null polynomial invariant for any given initial values if and only if there exists an eigenspace of $M$ with dimension at least 2 .
Example 4 We return to running Example 2. Matrix $M$ has 6 distinct eigenvalues, and so the corresponding eigenspaces are of dimension 1 . We denote by $E_{\lambda}$ the eigenspace corresponding to $\lambda$. Then $E_{4}$ has a basis $(1,-2,1,-6,6,9)^{\top}$, $E_{6}$ has a basis $(0,1,-1,2,-5,6)^{\top}, E_{9}$ has a basis $(0,0,1,0,4,4)^{\top}, E_{2}$ has a basis $(0,0,0,1,-1,-3)^{\top}, E_{3}$ has a basis $(0,0,0,0,1,2)^{\top}$, and $E_{1}$ has a basis $(0,0,0,0,0,1)^{\top}$. Also, suppose that the initiation step is given by $\left(x_{1}=0, x_{2}=-2\right)$, i.e., $\left(u_{1}, u_{2}\right)=(0,2)$, which corresponds to the hyperplane $Q(0,2)=0$ in $\mathbb{R}_{2}\left[x_{1}, x_{2}\right]$.

We start with simple initial conditions and consider general conditions in the sequel. Theorem 3 applies, and since it is clear that $(0,0,1,0,4,4)^{\top}$ belongs to the hyperplane, we get $x_{2}^{2}+4 x_{2}+4=0$ is an inductive invariant for that loop with these specific initial conditions.

Example 5 We study the following transition system [SSM04b], corresponding to the multiplication of 2 numbers, and where the transition is $\tau=\left\langle l_{i}, l_{i}, \rho_{\tau}\right\rangle$, with

$$
\rho_{\tau}=\left[\begin{array}{c}
s^{\prime}=s+i \\
j^{\prime}=j-1 \\
i^{\prime}=i \\
j_{0}^{\prime}=j_{0}
\end{array}\right]
$$

We need to find a $\lambda$ such that $Q\left(s+i, j+1, i, j_{0}\right)=\lambda Q\left(s, j, i, j_{0}\right)$.

- Step 1: We build the associated matrix $M$ :

$$
\left(\begin{array}{ccccccccccccccc}
1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
2 & 0 & 1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
1 & 0 & 1 & 0 & 0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 \\
0 & -1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & -2 & 0 & 0 & 0 & 0 & 0 & 0 & 1 & 0 & 0 & 0 \\
0 & -1 & 0 & 0 & 0 & -1 & 0 & 0 & 0 & 0 & 1 & 0 & 1 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & -1 & 0 & 0 & 0 & 0 & 0 & 0 & 1 & 0 \\
0 & 0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 & 0 & -1 & 0 & 0 & 1
\end{array}\right) .
$$

- Step 2: We compute the eigenvectors which will provide us with a basis for non-trivial $\lambda$-invariants. Here, an evident eigenvalue is 1.
- Step 3: It is clear, in view of the matrix $M$, that $\operatorname{dim}(\operatorname{Ker}(M-I)) \geq 2$. As the eigenspace associated to eigenvalue 1 is of dimension 2, Corollary 2 applies. For example, the vector

$$
(0,0,0,0,0,1,0,0,-1,0,1,0,0,0,0)^{\top}
$$

is the eigenvector corresponding to the $\lambda$-invariant $s+j i-i j_{0}$.
Note that without computing Gröbner bases or performing quantifier elimination, we found the invariant $s+j i-$ $i j_{0}=0$ obtained by Sankaranarayanan et al. in [SSM04b]. The consecution scale technique will give a non-null invariant whatever the initial values are, and this explains why a non-trivial invariant was found in that work.

### 4.3. Limits of constant-scale consecution

Here we consider an algebraic transition relation where the instructions are described by polynomials with degree greater than 1.

Example 6 Consider the following loop: $\rho_{\tau} \equiv\left[\begin{array}{c}x^{\prime}=x(y+1) \\ y^{\prime}=y^{2}\end{array}\right]$. At step $k$ of the iteration, this loop computes the sum $1+y+\ldots+y^{2^{k}-1}$. Let $P(x, y)=a_{0} x^{2}+a_{1} x y+a_{2} y^{2}+a_{3} x+a_{4} y+a_{5}$ be a candidate $\lambda$-invariant. With the Gröbner Bases $\left\{x^{\prime}-x(y+1), y^{\prime}-y^{2}\right\}$, and with the total-degree lexicographic ordering given by the precedence $x^{\prime}>y^{\prime}>x>y$, we can get the loop ideal of $\mathbb{K}\left[x^{\prime}, y^{\prime}, x, y\right]$. Modulo this loop ideal, we have $P\left(x^{\prime}, y^{\prime}\right)=P\left(x(y+1), y^{2}\right)$. Put $P^{\prime}(x, y)=P\left(x(y+1), y^{2}\right)$. After expanding we get $P^{\prime}(x, y)=a_{0} x^{2} y^{2}+a_{1} x y^{3}+$ $a_{2} y^{4}+2 a_{0} x^{2} y+a_{1} x y^{2}+a_{0} x^{2}+a_{3} x y+a_{4} y^{2}+a_{3} x+a_{5}$. If we try a constant-scale consecution with parameter $\lambda$ we obtain:

$$
\begin{aligned}
& a_{0}=0 \quad a_{1}=0 \quad a_{3}=\lambda a_{3} \\
& a_{1}=0 \quad a_{0}=\lambda a_{0} \quad \lambda a_{4}=0 \\
& a_{2}=0 \quad a_{3}=\lambda a_{1} \quad a_{5}=\lambda a_{5} \\
& 2 a_{0}=0 \quad a_{4}=\lambda a_{2} .
\end{aligned}
$$

After simplifications, we get $a_{0}=a_{1}=a_{2}=a_{3}=a_{4}=0$ and $a_{5}=\lambda a_{5}$. If $\lambda \neq 1$ then $a_{5}=0$, which leads to a null invariant. Otherwise, $\lambda=1$ and we obtain the constant invariant $a_{5}$. Also, the initial condition implies that the constant invariant $a_{5}$ is null. So, using a constraint-based approach with constant-scaling [SSM04b] we can obtain only constant or null, i.e. trivial, invariants.

In the following section, we show how we handle this problem.

## 5. Algebraic discrete transition systems

In this section, we approach non-linear discrete systems.

## 5.1. $T$-scale invariant generation

Consider an algebraic transition system: $\rho_{\tau} \equiv\left[\begin{array}{c}X_{1}^{\prime}=P_{1}\left(X_{1}, . ., X_{n}\right) \\ \vdots \\ X_{n}^{\prime}=P_{n}\left(X_{1}, . ., X_{n}\right)\end{array}\right]$, where the loop updates can be represented using polynomials in $\mathbb{R}\left[X_{1}, \ldots, X_{n}\right]$ of the forms $P_{i}\left(X_{1}, \ldots, X_{n}\right)=\sum_{i_{1} \ldots, i_{n}} a_{i_{1}, ., i_{n}} X_{1}{ }^{i_{1}} \ldots X_{n}{ }^{i_{n}}$, where the coefficients $a_{i_{1}, \ldots, i_{n}}$ are in $\mathbb{R}$. We have the following $T$-scale discrete invariant characterization.

Theorem 4 A polynomial $Q$ in $\mathbb{R}\left[X_{1}, . ., X_{n}\right]$ is a $T$-scale discrete invariant for polynomial-scale consecution with $a$ parametric polynomial $T \in \mathbb{R}\left[X_{1}, \ldots, X_{n}\right]$ for $\tau$ if and only if

$$
Q\left(P_{1}\left(X_{1}, . ., X_{n}\right), . ., P_{n}\left(X_{1}, . ., X_{n}\right)\right)=T\left(X_{1}, . ., X_{n}\right) Q\left(X_{1}, . ., X_{n}\right)
$$

Example 7 Reconsider Example 6. We now take $\left(y=y_{0}, x=1\right)$ as initial values. We want to obtain a polynomial scale consecution with a parametric polynomial $T(x, y)=b_{0} y^{2}+b_{1} x+b_{2} y+b_{3}$. We thus obtain $P^{\prime}(s, x)=$ $\left(b_{0} y^{2}+b_{1} x+b_{2} y+b_{3}\right) P(x, y)$. In other words, we obtain the following multi-parametric linear system with parameters $b_{0}, b_{1}, b_{2}, b_{3}$ :

$$
\begin{array}{rlrl}
a_{0} & =b_{0} a_{0} & 0 & =b_{2} a_{5}+b_{3} a_{4} \\
a_{1} & =b_{0} a_{1} & 0 & =b_{0} a_{4}+b_{2} a_{2} \\
a_{2} & =b_{0} a_{2} & a_{3} & =b_{1} a_{5}+b_{2} a_{3}+b_{3} a_{3} \\
a_{5} & =b_{3}=b_{0} a_{5}+b_{2} a_{4}+b_{3} a_{2} \\
0 & & a_{1}=a_{3} b_{0}+b_{1} a_{2}+b_{2} a_{1} \\
a_{0} & & =b_{1} a_{3}+b_{3} a_{0} &
\end{array}
$$

We now describe a decision procedure for choosing parameter values. Consider the first three equations and choose $b_{0}=1$. In this way we aim at a high degree invariant for, otherwise, the coefficients $a_{0}, a_{1}, a_{2}$ of the highest degree terms would be null. Then, we are lead to another system with $b_{1} a_{0}=0$. For the same reason, choose $b_{1}=0$. Then we have $b_{2} a_{0}=2 a_{0}$. As a direct consequence, $b_{2}$ is set to 2 . Since equation $b_{3} a_{0}=a_{0}$ is in the resulting system, $b_{3}$ is set to 1 . Finally, we obtain the following system:

$$
\begin{array}{r}
a_{3}+a_{1}=0 \\
a_{4}+2 a_{2}=0 \\
a_{2}-a_{5}=0
\end{array}
$$

Having less equations than variables, we will have a non-trivial solution for generating of $T$-invariants. Now, we consider the hyperplane corresponding to the initial values, that is, $a_{2} y_{0}^{2}+\left(a_{1}+a_{4}\right) y_{0}+a_{0}+a_{1}+a_{5}=0$. As there are six variables and four equations, we will have again a non-trivial solution. A possible solution is the vector $\left(y_{0}\left(1-y_{0}\right), 1,1,-1,-2,1\right)^{\top}$. So, $y_{0}\left(1-y_{0}\right) x^{2}+x y+y^{2}-x-2 y+1=0$ is an invariant. Note that $T(x, y)=y^{2}+y+1$.
Remark 1 That is a simple constraint-based procedure, which can fail in more complex cases. Shortly, we will present a superior technique, from a more encompassing point of view.

### 5.2. A general theory for discrete transitions and polynomial systems

If $Q \in \mathbb{R}\left[X_{1}, \ldots, X_{n}\right]$ is of degree $r$ and the maximal degree of the $P_{i}$ 's is $d$, then we are looking for a $T$ of degree $e=d r-r$. Write its ordered coefficients as $\lambda_{0}, \ldots, \lambda_{s}$, with $s+1$ being the number of monomials of degree inferior to $e$.

Let $M$ be the matrix, in the canonical basis of $\mathbb{R}_{r}\left[X_{1}, . ., X_{2}\right]$ and $\mathbb{R}_{d r}\left[X_{1}, . ., X_{n}\right]$, of the morphism

$$
\mathscr{M}:\left\{\begin{aligned}
\mathbb{R}_{r}\left[X_{1}, \ldots, X_{n}\right] & \rightarrow \mathbb{R}_{d r}\left[X_{1}, \ldots, X_{n}\right] \\
Q\left(X_{1}, \ldots, X_{n}\right) & \mapsto Q\left(P_{1}\left(X_{1}, . ., X_{n}\right), \ldots, P_{n}\left(X_{1}, . ., X_{n}\right)\right) .
\end{aligned}\right.
$$

Let $L$ be the matrix, in the canonical basis of $\mathbb{R}_{r}$ and $\mathbb{R}_{d r}$, of the morphism

$$
\mathscr{L}:\left\{\begin{aligned}
\mathbb{R}_{r}\left[X_{1}, \ldots, X_{n}\right] & \rightarrow \mathbb{R}_{d r}\left[X_{1}, \ldots, X_{n}\right] \\
P & \mapsto T P
\end{aligned}\right.
$$

Matrix $L$ has a very simple form: its non zero coefficients are the $\lambda_{i}$ 's, and it has a natural block decomposition. Now let $Q \in \mathbb{R}\left[X_{1}, \ldots, X_{n}\right]$ be a $T$-scale discrete invariant for a transition relation defined by the $P_{i}$ 's. Then,

$$
\begin{array}{rlrl}
Q\left(P_{1}\left(X_{1}, . ., X_{n}\right), \ldots, P_{n}\left(X_{1}, . ., X_{n}\right)\right) & =T\left(X_{1}, . ., X_{n}\right) Q\left(X_{1}, . ., X_{n}\right) & & \Leftrightarrow \\
\mathscr{M}(Q) & =\mathscr{L}(Q) & & \Leftrightarrow \\
(\mathscr{M}-\mathscr{L})(Q) & =0_{\mathbb{R}\left[X_{1}, \ldots, X_{n}\right]} & \Leftrightarrow \\
Q & \in \operatorname{Ker}(M-L) . & & \Leftrightarrow
\end{array}
$$

A $T$-scale discrete invariant is nothing else than a vector in the kernel of $M-L$. Our problem is equivalent to finding a $L$ such that $M-L$ has a non-trivial kernel.
Theorem 5 Consider $M$ as described above. Then, there will be a $T$-scale discrete invariant if and only if there exists a matrix L, corresponding to $P \mapsto T P$, such that $M-L$ has a nontrivial kernel. Further, any vector in the kernel of $M-L$ will give rise to a $T$-scale invariant.

We denote by $v(r)$ the dimension of $\mathbb{R}_{r}\left[X_{1}, \ldots, X_{n}\right]$. Again, the last column of $M$ is $(0, \ldots, 0,1)^{\top}$. The last column of $L$ is $\left(0, \ldots, 0, \lambda_{0}, \ldots, \lambda_{s}\right)^{\top}$. Hence, choosing every $\lambda_{i}$ to be zero, except for $\lambda_{s}=1$, the last column of $M-L$ will be null. With this choice of $L$ (or $T=1$ ), we get at least $T$-invariants corresponding to constant polynomials. Now, $M-L$ having a non-trivial kernel is equivalent to its rank being less than the dimension $v(r)$ of $\mathbb{R}_{r}\left[X_{1}, \ldots, X_{n}\right]$. This is equivalent to the fact that each $v(r) \times v(r)$ sub-determinant of $M-L$ is equal to zero [Lan02]. Those determinants are polynomials in variables $\left(\lambda_{0}, \lambda_{1}, \ldots, \lambda_{s}\right)$, which we will denote by $V_{1}\left(\lambda_{0}, \lambda_{1}, \ldots, \lambda_{s}\right), \ldots, V_{s}\left(\lambda_{0}, \lambda_{1}, \ldots, \lambda_{s}\right)$.
Theorem 6 There is a non-trivial T-scale invariant if and only if the polynomials $\left(V_{1}, . ., V_{s}\right)$ admit a common root, other than the trivial one $(0, \ldots, 0,1)$.
Remark 2 This theorem provides us with important existence results. But there is a more practical way of computing invariant ideals without computing common roots and sub-determinants. We will examine that in the next section.

We first study the general case of degree two algebraic transition systems with two variables in the loop. Such transition systems have the form: $\rho_{\tau} \equiv\left[\begin{array}{c}x^{\prime}=c_{0} x^{2}+c_{1} x y+c_{2} y^{2}+c_{3} x+c_{4} y+c_{5} \\ y^{\prime}=d_{0} x^{2}+d_{1} x y+d_{2} y^{2}+d_{3} x+d_{4} y+d_{5}\end{array}\right]$. In this case, matrices $M$ and $L$ will be as follows:

$$
\left.M=\left(\begin{array}{cccccc}
c_{0}^{2} & c_{0} d_{0} & d_{0}^{2} & 0 & 0 & 0 \\
2 c_{0} c_{1} & c_{0} d_{1}+c_{1} d_{0} & 2 d_{0} d_{1} & 0 & 0 & 0 \\
2 c_{0} c_{2}+c_{1}^{2} & c_{0} d_{2}+c_{1} d_{1}+c_{2} d_{0} & 2 d_{0} d_{2}+d_{1}^{2} & 0 & 0 & 0 \\
2 c_{1} d_{1} & c_{1} d_{2}+c_{2} d_{1} & 2 d_{1} d_{2} & 0 & 0 & 0 \\
c_{2}^{2} & c_{2} d_{2} & d_{2}^{2} & 0 & 0 & 0 \\
2 c_{0} c_{3} & c_{0} d_{3}+c_{3} d_{0} & 2 d_{0} d_{3} & 0 & 0 & 0 \\
2\left(c_{0} c_{4}+c_{1} c_{3}\right) & c_{0} d_{4}+c_{1} d_{3}+c_{3} d_{1}+c_{4} d_{0} & 2\left(d_{0} d_{4}+d_{1} d_{3}\right) & 0 & 0 & 0 \\
2\left(c_{1} c_{4}+c_{2} c_{3}\right) & c_{1} d_{4}+c_{2} d_{3}+c_{3} d_{2}+c_{4} d_{1} & 2\left(d_{1} d_{4}+d_{2} d_{3}\right) & 0 & 0 & 0 \\
2 c_{2} c_{4} & c_{2} d_{4}+c_{4} d_{2} & 2 d_{2} d_{4} & 0 & 0 & 0 \\
2 c_{0} c_{5}+c_{3}^{2} & c_{0} d_{5}+c_{3} d_{3}+c_{5} d_{0} & 2 d_{0} d_{5}+d_{3}^{2} & c_{0} & d_{0} & 0 \\
2\left(c_{1} c_{5}+c_{3} c_{4}\right) & c_{1} d_{5}+c_{3} d_{4}+c_{4} d_{3}+c_{5} d_{1} & 2\left(d_{1} d_{5}+d_{3} d_{4}\right) & c_{1} & d_{1} & 0 \\
2 c_{2} c_{5}+c_{4}^{2} & c_{2} d_{5}+c_{4} d_{4}+c_{5} d_{2} & 2 d_{2} d_{5}+d_{4}^{2} & c_{2} & d_{2} & 0 \\
2 c_{3} c_{5} & c_{3} d_{5}+c_{5} d_{3} & 2 d_{3} d_{5} & c_{3} & d_{3} & 0 \\
2 c_{4} c_{5} & c_{4} d_{5}+c_{5} d_{4} & 2 d_{4} d_{5} & c_{4} & d_{4} & 0 \\
c_{5}^{2} & c_{5} d_{5} & d_{5}^{2} & c_{5} & d_{5} & 1
\end{array}\right) L=\begin{array}{ccccc}
\lambda_{0} & 0 & 0 & 0 & 0 \\
\lambda_{1} & \lambda_{0} & 0 & 0 & 0 \\
\lambda_{2} & \lambda_{1} & \lambda_{0} & 0 & 0 \\
0 & \lambda_{2} & \lambda_{1} & 0 & 0 \\
0 \\
0 & 0 & \lambda_{2} & 0 & 0 \\
\lambda_{3} & 0 & 0 & \lambda_{0} & 0 \\
\lambda_{4} & \lambda_{3} & 0 & \lambda_{1} & \lambda_{0} \\
0 & 0 \\
\lambda_{4} & \lambda_{3} & \lambda_{2} & \lambda_{1} & 0 \\
0 & 0 & \lambda_{4} & 0 & \lambda_{2} \\
0 \\
\lambda_{5} & 0 & 0 & \lambda_{3} & 0 \\
\lambda_{0} \\
0 & \lambda_{5} & 0 & \lambda_{4} & \lambda_{3} \\
\lambda_{1} \\
0 & 0 & \lambda_{5} & 0 & \lambda_{4} \\
\lambda_{2} \\
0 & 0 & 0 & \lambda_{5} & 0 \\
\lambda_{3} \\
0 & 0 & 0 & 0 & \lambda_{5} \\
\lambda_{4} \\
0 & 0 & 0 & 0 & 0 \\
\lambda_{5}
\end{array}\right) .
$$

For the rank of $M-L$ to be less than 6 , one has to calculate each $6 \times 6$ sub-determinant obtained by canceling 9 lines of $M-L$. They will be polynomials of degree less than 6 in the variables $\left(\lambda_{0}, \ldots, \lambda_{5}\right)$. In this way $M-L$ will be of degree less than 6 if and only if $\left(\lambda_{0}, \ldots, \lambda_{5}\right)$ are roots of each of those polynomials. In many cases, it is easy to find a matrix $L$ such that $M-L$ has a non-trivial kernel. We describe and deal with several decidable classes (see Table 1b line 3 and the lines 13-20). The following important remark make clear the advances reached when comparing this approach to related constraint-based methods.
Remark 3 The unknown coefficients $\left(\lambda_{0}, \ldots, \lambda_{s}\right)$ appearing in the matrix $L$, used to approximate the consecution conditions, are assigned in order to insure that the nullspaces generated are not trivial ones. In fact, these unknown coefficients do not have the same roles as the templates used in constraint-based approaches. In our method, these parameters do not take part in a constraint solving problem. Instead, they allow us to obtain a sufficiently precise approximation to the consecution condition in order to guarantee the existence and the computation of vector spaces of $T$-invariants, that is, nullspaces.

### 5.3. Generating invariant ideals with an initiation step

Consider initial values given by unknown parameters $\left(X_{1}=u_{1}, \ldots, X_{n}=u_{n}\right)$. The initial step defines, on $\mathbb{R}_{r}\left[x_{1}, \ldots, x_{n}\right]$, a linear form $P \mapsto P\left(u_{1}, \ldots, u_{n}\right)$. Hence, initial values correspond to a hyperplane in $\mathbb{R}_{r}\left[X_{1}, . ., X_{n}\right]$, given by the kernel of $P \mapsto P\left(u_{1}, \ldots, u_{n}\right)$, which is $\left\{Q \in \mathbb{R}_{r}\left[X_{1}, . ., X_{n}\right] \mid Q\left(u_{1}, \ldots, u_{n}\right)=0\right\}$.
Theorem 7 Let $Q$ be in $\mathbb{R}_{r}\left[X_{1}, . ., X_{n}\right]$. Then $Q$ is an inductive invariant for the transition system with initial values $\left(u_{1}, . ., u_{n}\right)$ if and only if there exists a matrix $L \neq 0$, i.e, one of $P \mapsto T P$, corresponding to $T$ in $\mathbb{R}_{e}\left[X_{1}, . ., X_{n}\right]$, such that $Q$ is in the intersection of $\operatorname{Ker}(M-L)$ and the hyperplane given by the initial values $Q\left(u_{1}, \ldots, u_{n}\right)=0$. The invariants will correspond to vectors in the intersection.

Now, if $\operatorname{Dim}(\operatorname{Ker}(M-L)) \geq 2$ then $\operatorname{Ker}(M-L)$ will intersect any initial (semi-)hyperplane. We can state the following corollary, important in practice.
Corollary 3 There are non-trivial invariants for any given initial values if and only if there exists a matrix $L$ such that $\operatorname{Ker}(M-L)$ has dimension at least 2 . The basis of $\operatorname{Ker}(M-L)$ being a basis for non-trivial invariants.
There are non-trivial invariants for any given initial values if and only if there exists a matrix $L$, corresponding to multiplicative template in $T$, such that $\operatorname{Ker}(M-L)$ has dimension at least 2 .

### 5.4. Example

Consider the following transition: $\tau=\left\langle l_{i}, l_{j}, \rho_{\tau} \equiv\left[\begin{array}{c}x^{\prime}=x y+x \\ y^{\prime}=y^{2}\end{array}\right]\right\rangle$.

- Step 1: We build the matrix $M-L$. The maximal degree of $\rho_{\tau}$ is $d=2$, and so the $T$-scale invariant will be of degree $r=2$. Also, $T$ is of degree $e=d r-r=2$ and we write $\lambda_{0}, \ldots, \lambda_{5}$ as its ordered coefficients. Then its canonical form is $T=\lambda_{0} x^{2}+\lambda_{1} x y+\lambda_{2} y^{2}+\lambda_{3} x+\lambda_{4} y+\lambda_{5}$. Consider the associated morphisms $\mathscr{M}$ and $\mathscr{L}$ from $\mathbb{R}_{2}[x, y]$ to $\mathbb{R}_{4}[x, y]$. Using the basis $C_{1}=\left(x^{2}, x y, y^{2}, x, y, 1\right)$ of $\mathbb{R}_{2}[x, y]$ and the basis $C_{2}=\left(x^{4}, y x^{3}, y^{2} x^{2}, y^{3} x, y^{4}, x^{3}, x^{2} y, x y^{2}, y^{3}, x^{2}, x y, y^{2}, x, y, 1\right)$ of $\mathbb{R}_{4}[x, y]$, our algorithm compute the matrix $M-L$ as

$$
M-L_{\left(\lambda_{0}, \ldots, \lambda_{5}\right)}=\left(\begin{array}{cccccc}
-\lambda_{0} & 0 & 0 & 0 & 0 & 0 \\
-\lambda_{1} & -\lambda_{0} & 0 & 0 & 0 & 0 \\
1-\lambda_{2} & -\lambda_{1} & -\lambda_{0} & 0 & 0 & 0 \\
0 & 1-\lambda_{2} & -\lambda_{1} & 0 & 0 & 0 \\
0 & 0 & 1-\lambda_{2} & 0 & 0 & 0 \\
-\lambda_{3} & 0 & 0 & -\lambda_{0} & 0 & 0 \\
2-\lambda_{4} & -\lambda_{3} & 0 & -\lambda_{1} & -\lambda_{0} & 0 \\
0 & 1-\lambda_{4} & -\lambda_{3} & -\lambda_{2} & -\lambda_{1} & 0 \\
0 & 0 & -\lambda_{4} & 0 & -\lambda_{2} & 0 \\
1-\lambda_{5} & 0 & 0 & -\lambda_{3} & 0 & -\lambda_{0} \\
0 & -\lambda_{5} & 0 & 1-\lambda_{4} & -\lambda_{3} & -\lambda_{1} \\
0 & 0 & -\lambda_{5} & 0 & 1-\lambda_{4} & -\lambda_{2} \\
0 & 0 & 0 & 1-\lambda_{5} & 0 & -\lambda_{3} \\
0 & 0 & 0 & 0 & -\lambda_{5} & -\lambda_{4} \\
0 & 0 & 0 & 0 & 0 & 1-\lambda_{5}
\end{array}\right) .
$$

- Step 2: We now reduce the rank of $M-L$ by assigning values to the $\lambda_{i}$ 's. Our procedure fixes $\lambda_{0}=\lambda_{1}=\lambda_{3}=0$, $\lambda_{2}=\lambda_{5}=1$ and $\lambda_{4}=2$, so that $T(x, y)=y^{2}+2 y+1$. The first column of $M-L$ becomes zero and the second column is equal to the fourth. Hence, the rank of $M-L$ is less than 4 and its kernel has dimension at least 2 . Any vector in this kernel will be a $T$-invariant.
Before, proceeding to Step 3 we give more details on our rank reduction procedure which allows us to choose the coefficients $\lambda_{0}, \ldots, \lambda_{4}$ and $\lambda_{5}$ such that the matrix $M-L$ does not have a trivial kernel. Taking into consideration the specific type of matrix we are manipulating, we are going to determine for which values of $\lambda_{0}, \ldots, \lambda_{4}$ and $\lambda_{5}$, the rank of $M-L$ is minimal. We proceed by an analysis over the top non-zero elements of the columns of $M-L$ and the possible values that could be chosen for the parameters in order to decrease the actual rank. At each step of the assignments, we will echelon the matrix by making the highest term of one column to vanish. For that, we index the matrix $M-L$ as $M-L_{\left(\lambda_{0}, \ldots, \lambda_{5}\right)}$, in order to keep track of the assignment and rank obtained during the procedure.
Looking at the top non-zero elements of $M-L_{\left(\lambda_{0}, \ldots, \lambda_{5}\right)}$, consider first the parameter $\lambda_{0}$. If $\lambda_{0}$ is not zero, we obtain an echelon form matrix of rank 6 , which is maximal. Thus $\lambda_{0}$ is fixed to 0 . Now, considering the top non-zero elements of the columns of $M-L_{\left(0, \lambda_{1}, \ldots, \lambda_{5}\right)}$, the procedure fixes $\lambda_{1}$ to 0 for the same reason. We obtain the following matrix:

$$
M-L_{\left(0,0, \lambda_{2}, \ldots, \lambda_{5}\right)}=\left(\begin{array}{cccccc}
0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 \\
1-\lambda_{2} & 0 & 0 & 0 & 0 & 0 \\
0 & 1-\lambda_{2} & 0 & 0 & 0 & 0 \\
0 & 0 & 1-\lambda_{2} & 0 & 0 & 0 \\
-\lambda_{3} & 0 & 0 & 0 & 0 & 0 \\
2-\lambda_{4} & -\lambda_{3} & 0 & 0 & 0 & 0 \\
0 & 1-\lambda_{4} & -\lambda_{3} & -\lambda_{2} & 0 & 0 \\
0 & 0 & -\lambda_{4} & 0 & -\lambda_{2} & 0 \\
1-\lambda_{5} & 0 & 0 & -\lambda_{3} & 0 & 0 \\
0 & -\lambda_{5} & 0 & 1-\lambda_{4} & -\lambda_{3} & 0 \\
0 & 0 & -\lambda_{5} & 0 & 1-\lambda_{4} & -\lambda_{2} \\
0 & 0 & 0 & 1-\lambda_{5} & 0 & -\lambda_{3} \\
0 & 0 & 0 & 0 & -\lambda_{5} & -\lambda_{4} \\
0 & 0 & 0 & 0 & 0 & 1-\lambda_{5}
\end{array}\right) .
$$

Again, we look at the non-zero top column elements with $\lambda_{2}$. If $\lambda_{2}$ is not in $\{1,0\}$, we obtain an echelon form matrix of rank 6 , which is maximal. Thus we need to consider the case where $\lambda_{2}$ is assigned a value in $\{1,0\}$. Hence $\lambda_{2}$ is first assigned to 1 . Then $\lambda_{3}$ has to be 0 , otherwise the rank will be maximal. The procedure is now working with:

$$
M-L_{\left(0,0,1,0, \lambda_{4}, \lambda_{5}\right)}=\left(\begin{array}{cccccc}
0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 \\
2-\lambda_{4} & 0 & 0 & 0 & 0 & 0 \\
0 & 1-\lambda_{4} & 0 & -1 & 0 & 0 \\
0 & 0 & -\lambda_{4} & 0 & -1 & 0 \\
1-\lambda_{5} & 0 & 0 & 0 & 0 & 0 \\
0 & -\lambda_{5} & 0 & 1-\lambda_{4} & 0 & 0 \\
0 & 0 & -\lambda_{5} & 0 & 1-\lambda_{4} & -1 \\
0 & 0 & 0 & 1-\lambda_{5} & 0 & 0 \\
0 & 0 & 0 & 0 & -\lambda_{5} & -\lambda_{4} \\
0 & 0 & 0 & 0 & 0 & 1-\lambda_{5}
\end{array}\right) .
$$

In order to reduce its rank by making the highest term of a column vanish, $\lambda_{4}$ is assigned a value in $\{2,1,0\}$. Then, $\lambda_{4}$ is first fixed to 2 which implies that $\lambda_{5}$ needs to be assigned to 1 . We conclude that with $\left(\lambda_{0}, \ldots, \lambda_{4}, \lambda_{5}\right)=(0,0,1,0,2,1)$ the matrix $M-L$ does not have a trivial kernel and $T(x, y)=y^{2}+2 y+1$ can be used to generate a vector space of $T$-invariants. The procedure does not stop here. It continues to consider the other possibilities for $\lambda_{2}, \lambda_{4}$ and $\lambda_{5}$, thus generating more polynomials for approximating scaling consecution, leading to other vector spaces of polynomial scale and inductive invariants. With $\lambda_{2}=1$ and $\lambda_{4}=1$ or $\lambda_{4}=0$ one has to fix $\lambda_{5}=0$ and generate the polynomials $T_{2}(x, y)=y^{2}+y$ and $T_{3}(x, y)=y^{2}$. It remains to treat the case $\lambda_{2}=0$ and the possibilities thereof. One needs to consider again $M-L_{\left(0,0, \lambda_{2}, \ldots, \lambda_{5}\right)}$. In this case, one needs again to fix $\lambda_{3}$ to 0 . We obtain the matrix:

$$
M-L_{\left(0,0,0,0, \lambda_{4}, \lambda_{5}\right)}=\left(\begin{array}{cccccc}
0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 \\
1 & 0 & 0 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 & 0 & 0 \\
0 & 0 & 1 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 \\
2-\lambda_{4} & 0 & 0 & 0 & 0 & 0 \\
0 & 1-\lambda_{4} & 0 & 0 & 0 & 0 \\
0 & 0 & -\lambda_{4} & 0 & 0 & 0 \\
1-\lambda_{5} & 0 & 0 & 0 & 0 & 0 \\
0 & -\lambda_{5} & 0 & 1-\lambda_{4} & 0 & 0 \\
0 & 0 & -\lambda_{5} & 0 & 1-\lambda_{4} & 0 \\
0 & 0 & 0 & 1-\lambda_{5} & 0 & 0 \\
0 & 0 & 0 & 0 & -\lambda_{5} & -\lambda_{4} \\
0 & 0 & 0 & 0 & 0 & 1-\lambda_{5}
\end{array}\right) .
$$

Then, looking at the top column elements of $M-L_{\left(0,0,0,0, \lambda_{4}, \ldots, \lambda_{5}\right)}$, we only have now two possibilities for $\lambda_{4}$ : it has to be assigned a value in $\{1,0\}$. When $\lambda_{4}=1$ we need to treat $\lambda_{5}=1$ and $\lambda_{5}=0$ leading to the polynomials $T_{4}(x, y)=y+1$ and $T_{5}(x, y)=y$. Now, treating the other case, with $\lambda_{4}=0$, one has no choice but to assign $\lambda_{5}$ to 1 leading to the polynomial $T_{6}(x, y)=1$.
Conclusion of Step 2: The polynomials $T, T_{2}, \ldots, T_{6}$ would guarantee a non-trivial Kernel for $M-L$. We move to Step 3 considering $T(x, y)=y^{2}+2 y+1$ because with $\left(\lambda_{0}, \ldots, \lambda_{5}\right)$, the rank of $M-L$ is minimal (i.e., null space of dimensions 3 ).

- Step 3: Now matrix $M-L$ satisfies the hypotheses of Theorem 5. So, there will always be invariants, whatever the initial values. We compute a basis of $\operatorname{Ker}(M-L):[(1,0,0,0,0,0),(0,1,0,-1,0,0),(0,0,1,0,-2,1)]$. The vectors of the basis are interpreted in the canonical basis $C_{1}$ of $\mathbb{R}_{2}[x, y]$, giving: $\left\{\left\{x^{2}, x y-x, y^{2}-2 y+1\right\}\right\}$

We have obtained an ideal for non-trivial inductive invariants. In other words, for all $G_{1}, G_{2}, G_{3} \in \mathbb{R}[x, y]$, $G_{1}(x, y)\left(x^{2}\right)+G_{2}(x, y)(x y-x)+G_{3}(x, y)\left(y^{2}-2 y+1\right)=0$ is an inductive invariant. For instance, consider the initial step $\left(y=y_{0}, x=1\right)$. A possible invariant is $y_{0}\left(1-y_{0}\right) x^{2}+x y-x+y^{2}-2 y+1=0$. We can also consider $T_{2}(x, y)=y^{2}+y$ and $T_{4}(x, y)=y+1$ as they both provide kernels of dimension at least 2 , leading to the two other vector spaces of non-trivial inductive invariants.

## 6. Obtaining optimal degree bounds for discrete transition systems

In order to guarantee the existence of non-trivial invariants, we look for a polynomial $T$ such that $\operatorname{Ker}(M-L) \neq 0$. The pseudo code depicted in Algorithm 1 illustrates the strategy. Its contribution relies on very general sufficient conditions for the existence and computation of invariants.

As input we have $r$, the candidate degree for the basis invariant elements, and $P_{1}, . . P_{n}$, the $n$ polynomials given by the transition relation in the loop program. We first compute $d$, the maximal degree of the $P_{i}$ 's as can be seen by Max_degree $\left(\left\{P_{1}, \ldots, P_{n}\right\}\right)$, at line 4 . Following the instructions provided in Sect. 4 the function Matrix_D $\left(r, d r, P_{1}, \ldots, P_{n}\right)$ construct the matrix of the morphism $\mathscr{M}$. Then, we detail the cases were the transitions are defined by non-linear systems, i.e., when $d \leq 2$. See the condition at line 7 . Then, we define $T$ as a polynomial of degree $d r-r$ in its canonical form, i.e., with parameterized coefficients. See Template_Canonical_Form ( $n, d, r, d r-r$ ), at line 7. Here, Template_Canonical_Form $(n, d, r, d r-r)$ returns the lists of unknown coefficient that we denoted by $\lambda_{0}, \ldots, \lambda_{s}$ in Sect. 5.2. Note that $s$ depends only on the values of $n, d$, and $r$. One can now call the function Matrix_L $(r, d r, T)$ that construct the matrix $L$ as shown in Sect. 5.2.

Next, we apply our decision procedure Reduce_Rank_Assigning_Values $(M-L)$ to assign values to the coefficients of $T$ in such a way that $\operatorname{Ker}(M-L) \neq 0$. See line 10. As we saw in the previous section, $\operatorname{Ker}(M-L) \neq 0$ is equivalent to having $\operatorname{Rank}(M-L)<\operatorname{Dim}\left(\mathbb{R}_{r}\left[X_{1}, \ldots, X_{n}\right]\right)$. In other words, it is the same as having $M-L$ with rank strictly less than the dimension $v(r)$ of $\mathbb{R}_{r}\left[X_{1}, \ldots, X_{n}\right]$. We then reduce the rank of $M-L$ by assigning values of the parameters in $L$. This function applies the decision procedure detailed in Sect. 5.4, Step 2. Next, we determine whether the matrix obtained, $\overline{M-L}$, has a trivial kernel by first computing its rank and then checking if $\operatorname{Rank}(\overline{M-L})<\operatorname{Dim}\left(R_{r}\left[X_{1}, . ., X_{n}\right]\right)$ holds, at line 11. We can now apply our main Theorem 5:

- If $\overline{M-L}$ has a trivial kernel, we know there is no $T$-scale invariants of degree less than $r$ and we can increase the degree $r$ of the desired invariants until Theorem 5, or Corollary 3, applies, or until stronger hypotheses occur, e.g. if all $v(r) \times v(r)$ sub-determinants are null. Note, at line 12, the call to return Ideal_Loop_Inv_Gen $(r+$ $1, P_{1}, \ldots, P_{n}, X_{1}, \ldots, X_{n}$ ). If there is no ideal for non-trivial invariants for a value $r_{i}$ then we conclude that there is no ideal of non-trivial invariants for all degrees $k \leq r_{i}$. This can also be used to guide other constraint-based techniques, since checking for invariance with a template of degree less or equal to $r_{i}$ will not be necessary.
- Otherwise, Theorem 5, or Corollary 3, guarantees the existence and computation of $T$-invariants. Finally, the function Nullspace_Basis $(\overline{M-L})$ outputs the basis of the nullspace of the matrix $\overline{M-L}$, in order to construct non-trivial invariants. See line 15.

For basis computations, we use well-known state-of-the-art algorithms, for example those that Sage provides. These algorithms calculate the eigenvalues and associated eigenspaces of $\overline{M-L}$ when it is a square matrix. When $\overline{M-L}$ is a rectangular matrix, we can use its singular value decomposition (SVD). A SVD of $\overline{M-L}$ provides an explicit representation of its rank and kernel by computing unitary matrices $U$ and $V$ and a regular diagonal matrix $S$ such that $\overline{M-L}=U S V$. We compute the SVD of a $v(r+d-1) \times v(r)$ matrix $\bar{M}$ in two steps. First, we reduce it to a bi-diagonal matrix, with a cost of $O\left(v(r)^{2} v(r+d-1)\right)$ flops. The second step relies on an iterative method, as is also the case for other algorithms that compute eigenvalues. In practice, however, it suffices to compute the SVD up to a certain precision, i.e. up to a machine epsilon. In this case, the second step takes $O(v(r))$ iterations, each using $O(v(r))$ flops. So, the overall cost is $O\left(v(r)^{2} v(r+d-1)\right)$ flops. For the encoding of the algorithm we could rewrite Corollary 3 as follows.
Corollary 4 Let $(M-L)^{\prime}=U S V$ be the singular value decomposition of matrix $(M-L)^{\prime}$ described just above. There will be a non-trivial T-invariant for any given initial condition if and only if the number of non-zero elements in matrix $S$ is less than $v(r)-2$, where $v(r)$ is the dimension of $\mathbb{R}_{r}\left[x_{1}, \ldots, x_{n}\right]$. Moreover, the orthonormal basis for the nullspace obtained from the decomposition directly gives an ideal for non-linear invariants.

```
Algorithm 1: Ideal Loop_Inv_Gen \(\left(r, P_{1}, \ldots, P_{n}, X_{1}, \ldots, X_{n}\right)\)
    /*Finding degree bounds for discrete transitions. \({ }^{*}\);
    Data: \(r\) is the candidate degree for the set of basis invariants elements we are looking for, \(P_{1}, . . P_{n}\) the \(n\)
                polynomials given by the considered loop, and \(X_{1}, . . X_{n} \in V\)
    Result: Ideal_Inv, a basis of ideal of invariants.
    begin
        int \(d\);
        Template \(T\);
        Matrix \(M, L\);
        \(d \longleftarrow\) Max_degree \(\left(\left\{P_{1}, \ldots, P_{n}\right\}\right) ;\)
        \(/ * d\) is the maximal degree of \(P_{i}\) 's \({ }^{*} /\);
        \(M \longleftarrow\) Matrix_D \(\left(r, d r, P_{1}, \ldots, P_{n}\right)\);
        if \(d>=2\) then
            \(T \longleftarrow\) Template_Canonical_Form \((n, d, r, d r-r)\);
            \(L \longleftarrow \operatorname{Matrix} \_\mathbf{L}(r, d r, T)\);
            \(\overline{M-L} \longleftarrow\) Reduce_Rank_Assigning_Values \((M-L)\);
            if \(\operatorname{Rank}(\overline{M-L})>=\operatorname{Dim}\left(R_{r}\left[X_{1}, . ., X_{n}\right]\right)\) then
                    return Ideal_Loop_Inv_Gen \(\left(r+1, P_{1}, \ldots, P_{n}, X_{1}, \ldots, X_{n}\right)\);
                    /*We need to increase the degree r of candidates invariants.*/;
            else
                    return Nullspace_Basis( \(\overline{M-L}\) );
                    /*There exists an ideal of invariants that we can compute*/;
        else
            ... /*We refer to our previous work for constant scaling.*/;
    end
```

Remark 4 It is important to emphasize that eigenvectors or nullspace of $\overline{M-L}$ are computed after the parameters of $L_{T}$ have been assigned. When the discrete transition system has several variables and none or few parameters, which correspond to practical cases, $\overline{M-L}$ will be over the reals and there will be no need to use the symbolic version of these algorithms.

## 7. Invariant generation for discrete transitions and fractional systems

We now want to deal with transition systems $\rho_{\tau}$ of the following type

$$
\left[\begin{array}{c}
X_{1}^{\prime}=\frac{P_{1}\left(X_{1}, \ldots, X_{n}\right)}{Q_{1}\left(X_{1}, \ldots, X_{n}\right)} \\
\vdots \\
X_{n}^{\prime}=\frac{P_{n}\left(X_{1}, \ldots, X_{n}\right)}{Q_{n}\left(X_{1}, \ldots, X_{n}\right)},
\end{array}\right]
$$

where the $P_{i}$ 's and $Q_{i}$ 's belong to $\mathbb{R}\left[X_{1}, \ldots, X_{n}\right]$ and each $P_{i}$ is relatively prime to the corresponding $Q_{i}$. In this case, one needs to relax the consecution conditions to fractional-scale as soon as fractions appear in the transition relation.

Theorem 8 A polynomial $Q$ in $\mathbb{R}\left[X_{1}, \ldots, X_{n}\right]$ is a $F$-scale invariant for fractional discrete scale consecution with a parametric fractional $F \in \mathbb{R}\left(X_{1}, \ldots, X_{n}\right)$ for $\tau$ if and only if

$$
Q\left(\frac{P_{1}}{Q_{1}}, \ldots, \frac{P_{n}}{Q_{n}}\right)=F Q .
$$

Let $d$ be the maximal degree of the $P_{i}$ 's and $Q_{i}$ 's, and let $\Pi$ be the least common multiple of the $Q_{i}$ 's. Now let $U=X_{1}{ }^{i_{1}} . . X_{n}{ }^{i_{n}}$ be a monomial of degree less than $r$, i.e., $i_{1}+. .+i_{n}<r$. Then,

$$
\Pi^{r} U\left(P_{1} / Q_{1}, \ldots, P_{n} / Q_{n}\right)=\Pi^{r}\left(P_{1} / Q_{1}\right)^{i_{1}} \ldots\left(P_{n} / Q_{n}\right)^{i_{n}}
$$

But as $Q_{j}^{i_{j}}$ divides $\Pi^{i_{j}}$, for all $j$, we see that $Q_{1}^{i_{1}} \ldots Q_{n}^{i_{n}}$ divides $\Pi^{i_{1}+\ldots+i_{r}}$, which divides $\Pi^{r}$. We conclude that $\Pi^{r} Q\left(P_{1} / Q_{1}, \ldots, P_{n} / Q_{n}\right)$ is a polynomial for every $Q$ in $\mathbb{R}_{r}\left[X_{1}, . ., X_{n}\right]$.

Now suppose that $F=T / S$, with $T$ relatively prime to $S$, satisfies the equality of the previous theorem. Suppose, further, that we are looking for bases for invariants $Q$ of degree $r$. Then, multiplying by $\Pi^{r}$ we get

$$
\Pi^{r} Q\left(P_{1} / Q_{1}, \ldots, P_{n} / Q_{n}\right)=\left(\Pi^{r} T Q\right) / S
$$

As we have no a priory information on $Q$, in most cases $Q$ will be relatively prime to $S$. In this situation we see that $S$ divides $\Pi^{r}$. So, let $F$ be of the form $T / \Pi^{r}$, and note that we argued that this constraint is weak.

Now let $\mathscr{M}$ be the morphism

$$
\mathscr{M}:\left\{\begin{aligned}
\mathbb{R}_{r}\left[X_{1}, \ldots, X_{n}\right] & \rightarrow \mathbb{R}_{n r d}\left[X_{1}, \ldots, X_{n}\right] \\
Q & \mapsto \Pi^{r} Q\left(\frac{P_{1}}{Q_{1}}, \ldots, \frac{P_{n}}{Q_{n}}\right) .
\end{aligned}\right.
$$

Let $M$ be its matrix representation in the canonical basis, let $T$ be a polynomial in $\mathbb{R}_{n r d-r}\left[X_{1}, . ., X_{n}\right]$, and let $\mathscr{L}$ denote the vector space morphism

$$
\mathscr{L}:\left\{\begin{aligned}
\mathbb{R}_{r}\left[X_{1}, \ldots, X_{n}\right] & \rightarrow \mathbb{R}_{n r d}\left[X_{1}, \ldots, X_{n}\right] \\
Q & \mapsto T Q
\end{aligned}\right.
$$

Also, let $L$ be its matrix representation in the canonical basis. As stated in the following theorem, our problem is equivalent to finding a $L$ such that $M-L$ has a non-trivial kernel.
Theorem 9 Consider $M$ and $L$ as described above. Then, there exist $F$-scale invariants, where $F$ is of the form $T / \Pi^{r}$, if and only if there exists a matrix $L$ such that $\operatorname{Ker}(M-L) \neq \emptyset$. In this situation, any vector in the kernel of $M-L$ will give rise to a $F$-scale discrete invariant.
This is similar to Theorems 6 and 7. For the initiation step, we have a hyperplane in $\mathbb{R}_{r}\left[X_{1}, \ldots, X_{n}\right]$. In order for the transition system to make sense, the $n$-tuple of initial values must not be a root of any of the $Q_{i}$ 's, and similarly for further iterations as long as the loop is applied. In this way, they will not cancel $\Pi^{r}$. We have the following result.
Theorem 10 We have a non-trivial invariant if and only if there exists a matrix $L$ such that the intersection of the kernel of $M-L$ and the hyperplane given by the initial values is not zero. The invariants will correspond to vectors in the intersection.

We also have the following important corollary.
Corollary 5 We will have a non-trivial invariant for any non-trivial initial value if and only if there exists a matrix $L$ such that the dimension of $\operatorname{Ker}(M-L)$ is at least 2 .
Example 8 Consider the system $\rho_{\tau} \equiv\left[\begin{array}{c}x_{1}^{\prime}=\frac{x_{2}}{\left(x_{1}+x_{2}\right)} \\ x_{2}^{\prime}=\frac{x_{1}}{\left(x_{1}+2 x_{2}\right)}\end{array}\right]$. We are looking for $F$-scale invariant polynomials of degree 2. The least common multiple of $\left(x_{1}+x_{2}\right)$ and $\left(x_{1}+2 x_{2}\right)$ is their product, so that $\mathscr{M}$ is given by:

$$
Q \in \mathbb{R}_{2}\left[x_{1}, x_{2}\right] \mapsto\left[\left[\left(x_{1}+x_{2}\right)\left(x_{1}+2 x_{2}\right)\right]^{2} Q\left(\frac{x_{1}}{\left(x_{1}+x_{2}\right)}, \frac{x_{2}}{\left(x_{1}+2 x_{2}\right)}\right)\right] .
$$

As both $\frac{x_{2}}{\left(x_{1}+x_{2}\right)}$ and $\frac{x_{1}}{\left(x_{1}+2 x_{2}\right)}$ have degree zero,
$\left[\left(x_{1}+x_{2}\right)\left(x_{1}+2 x_{2}\right)\right]^{2} Q\left(\frac{x_{2}}{\left(x_{1}+x_{2}\right)}, \frac{x_{1}}{\left(x_{1}+2 x_{2}\right)}\right)$
will be a linear combination of degree 4 , if it is non-null.
Hence, $\mathscr{M}$ has values in $\operatorname{Vect}\left(x_{1}^{4}, x_{1}^{3} x_{2}, x_{1}^{2} x_{2}^{2}, x_{1} x_{2}^{3}, x_{2}^{4}\right)$. With $T$ and $Q$ in $\mathbb{R}_{2}\left[x_{1}, x_{2}\right]$ we verify that

$$
\left[\left(x_{1}+x_{2}\right)\left(x_{1}+2 x_{2}\right)\right]^{2} Q\left(\frac{x_{2}}{\left(x_{1}+x_{2}\right)}, \frac{x_{1}}{\left(x_{1}+2 x_{2}\right)}\right)=T Q
$$

As the left member is in $\operatorname{Vect}\left(x_{1}^{4}, x_{1}^{3} x_{2}, x_{1}^{2} x_{2}^{2}, x_{1} x_{2}^{3}, x_{2}^{4}\right), T$ must be of the form $\lambda_{0} x_{1}^{2}+\lambda_{1} x_{1} x_{2}+\lambda_{2} x_{2}^{2}$, and $Q$ must be of the form $a_{0} x_{1}^{2}+a_{1} x_{1} x_{2}+a_{3} x_{2}^{2}$. We see that we can take $Q$ in $\operatorname{Vect}\left(x_{1}^{2}, x_{1} x_{2}, x_{2}^{2}\right)$, and similarly for $T$. Then both $\mathscr{M}, \mathscr{L}: Q \mapsto T Q$ are morphisms from $\operatorname{Vect}\left(x_{1}^{2}, x_{1} x_{2}, x_{2}^{2}\right)$ into $\operatorname{Vect}\left(x_{1}^{4}, x_{1}^{3} x_{2}, x_{1}^{2} x_{2}^{2}, x_{1} x_{2}^{3}, x_{2}^{4}\right)$. In the corresponding canonical basis, the matrix $M-L$ is

$$
M-L=\left(\begin{array}{ccc}
-\lambda_{0} & 0 & 1 \\
-\lambda_{1} & 1-\lambda_{0} & 2 \\
1-\lambda_{2} & 3-\lambda_{1} & 1-\lambda_{0} \\
4 & 2-\lambda_{2} & -\lambda_{1} \\
4 & 0 & -\lambda_{2}
\end{array}\right)
$$

Taking $\lambda_{0}=1, \lambda_{1}=3$ and $\lambda_{2}=2$, the second column cancels out and the kernel will be equal to $\operatorname{Vect}(0,1,0)$. Now, Corollary 5 applies to $M-L$, and we obtain: $T\left(x_{1}, x_{2}\right) / Q\left(x_{1}, 2\right)=1 /\left(\left(x_{1}+x_{2}\right)\left(x_{1}+2 x_{2}\right)\right)^{2}$ and we have the nullspace $[(0,1,0)]$ and the basis of scale invariant $\left\{x_{1} x_{2}\right\}$. It was clear from the beginning that the corresponding polynomial $x_{1} x_{2}$ is $\frac{1}{\left[\left(x_{1}+x_{2}\right)\left(x_{1}+2 x_{2}\right)\right]^{2}}$-scale invariant. In particular, it is an invariant for the initial values $(0,1)$. Moreover, it clearly never cancels $x_{1}+x_{2}$ and $x_{1}+2 x_{2}$, because they are of the form $(a, 0)$ or $(0, b)$ with $a$ and $b$ strictly positive.

## 8. Branching conditions and nested loops

We have generated bases of vector spaces describing invariants for transition systems. A global invariant would be any invariant which is in the intersection of these vector spaces. In this way, we avoid the definition of a single isomorphism for the whole transition system. Instead, we generate the basis for each separate consecution condition. To compute a basis of global invariants, we could use the following theorem. It suggests to multiply all the elements of each computed basis. By so doing, we also avoid the heavy computation of ideal intersections.

Theorem 11 Let $I=\left\{I_{1}, \ldots, I_{k}\right\}$ a set of ideals in $\mathbb{R}\left[X_{1}, \ldots, X_{n}\right]$ such that $I_{j}=\left(f^{(j)}{ }_{1}, \ldots, f_{n_{j}}^{(j)}\right)$ for $j \in[1$, $k]$. Let $\otimes\left(I_{1}, \ldots, I_{k}\right)=\left\{\delta_{1}, \ldots, \delta_{n_{1} n_{2} \ldots n_{k}}\right\}$ be such that all elements $\delta_{i}$ in $\otimes\left(I_{1}, \ldots, I_{k}\right)$ are formed by the product of one element from each ideal in I. Assume that all $I_{j}$ 's are ideals for invariants for a loop at location $l_{j}$, described by a transition $\tau_{j}$. If all $l_{j}$ describe the same location or program point $l$, then $\otimes\left(I_{1}, \ldots, I_{k}\right)$ is an ideal of non-trivial non-linear invariants for the entire loop located at $l$.

Note that when we have several transitions looping at the same point, we can obtain an encoding of possible execution paths of a loop containing conditional statements.

This approach is a sound, but not complete, way of computing ideals for global invariants, and it also has a low computational time complexity. In order to take into account initial conditions we intersect these vector spaces with the initial semi-hyperplanes deduced from the isomorphism associated with initial requirements. Next, we show how our method deals with the conditional statements inside loops. Let's consider the following type of loop while (B_1) \{ [I_1;] if (B_2) $\left\{\left[I_{\_} 2 ;\right]\right\}$ else $\left\{\left[I_{-} ;\right]\right\}$[I_4;], where each $I_{i}$ represent a block of multivariate fractional instructions. First we represent the loop with the following two transitions $\tau_{1}=\left\langle l_{i}, l_{i},\left(\mathcal{B}_{1} \wedge \mathcal{B}_{2}\right), \rho_{\tau_{1}}\right\rangle$ and $\tau_{2}=\left\langle l_{i}, l_{i},\left(\mathcal{B}_{1} \wedge \neg \mathcal{B}_{2}\right), \rho_{\tau_{2}}\right\rangle$, where: $\rho_{\tau_{1}} \equiv\left[x_{1}^{\prime}=F_{1,\left[I_{1} ; I_{2} ; I_{4} ;\right]_{0}\left(x_{1}, \ldots x_{n}\right), \ldots, x_{n}^{\prime}=}=\right.$ $\left.F_{n,\left[I_{1} ; I_{2} ; I_{4} ;\right]_{0}}\left(x_{1}, \ldots, x_{n}\right)\right]$ and $\rho_{\tau_{2}} \equiv\left[x_{1}^{\prime}=F_{1,\left[I_{1} ; I_{3} ; I_{4} ;\right]_{0}}\left(x_{1}, \ldots x_{n}\right) \ldots, x_{n}^{\prime}=F_{n,\left[I_{1} ; I_{3} ; I_{4} ;\right]_{0}}\left(x_{1}, \ldots, x_{n}\right)\right]$, with $[; ;]$ 。 denoting our operator, based on separation rewriting rules, used to compose blocks of instructions. We first independently generate the ideals of invariants $\xi_{1}=\left(\mu_{1}, \ldots, \mu_{n}\right)$ and $\xi_{2}=\left(\kappa_{1}, \ldots, \kappa_{p}\right)$ for the respective transitions $\tau_{1}$ and $\tau_{2}$. Any element $\mu_{i} \in \xi_{1}$ refers to an inductive invariant $\mu_{i}\left(X_{1}, \ldots, X_{n}\right)=0$ corresponding to the partial loops described by transition $\tau_{1}$. Similarly, any $\kappa_{i} \in \xi_{2}$ refers to an inductive invariant $\kappa_{i}\left(X_{1}, \ldots, X_{n}\right)=0$ for the loop described by transition $\tau_{2}$. Then we can take $\mu_{i}\left(X_{1}, \ldots, X_{n}\right) * \kappa_{i}\left(X_{1}, \ldots, X_{n}\right)=0$ as global loop invariant, since these invariants will remain true in any sequence of transitions during the execution of the loop. We deal with loop conditions using the same methods that we proposed to handle initiation conditions. We know, for instance, that if our Corollary 3 holds, then there exist invariants for any (semi-)hyperplane that could be induced by the loop conditions. We illustrate this point in Fig. 1. Let $\left(P_{i}\left(x_{1}, . ., x_{n}\right)<0\right)$ be semi-algebraic loop conditions at location $l$ and let $Q$ be an inductive invariant for $\mathcal{D}(l)$. Thus $\left(P_{i}\left(x_{1}, . ., x_{n}\right)-Q\left(x_{1}, . ., x_{n}\right)<0\right)$ is also an inductive invariant. Then, we can build an operator, similar to the one introduced in Theorem 11, to generate, in a different way, ideals of non-trivial invariants at a state $l$ with semi-algebraic loop conditions. If a loop condition has the form $C_{i}\left(x_{1}, . ., x_{n}\right)=0$ we could then associate it directly to polynomial systems induced by the transition relations.


Fig. 1. Intersection between the conditional loop: $800<(x-5)^{2}+(y-5)^{2}+\left(y_{0}-5\right)^{2}<1000$ and the invariant $y_{0}\left(1-y_{0}\right) x^{2}+x y-x+y^{2}-2 y+1=$ 0 from the invariant ideal $\left(\left\{x^{2}, x y-x, y^{2}-2 y+1\right\}\right)$ obtained for Example 7

Example 9 Consider the following loop.

```
int u_0; //initialization
\(\left((M>0) \& \&(Z=1) \& \&\left(U=u_{-} 0\right) \ldots\right)\)
    While ((X>=1) || (Z>=z_0))\{
        If \((\mathrm{Y}>\mathrm{M})\) \{
            \(X=Y /(X+Y) ;\)
            \(\mathrm{Y}=\mathrm{X} /(\mathrm{X}+2 * \mathrm{Y}) ;\}\)
        Else\{
            \(Z=Z *(U+1) ;\)
            \(\left.U=U^{\wedge} 2 ;\right\}\)
    \}
```

We first generate an invariant for the loop corresponding to the first conditional if. Using Fractional-Scaling we obtain the basis of scale invariant $\{\{x y\}\}$. See Example 8 for more details. Then, we obtain the basis of invariants $\left\{u_{0} z^{2}-u_{0}^{2} z^{2}+z u+u^{2}-z-2 u+1, \ldots\right\}$ corresponding to the other alternative transition $\tau_{2}$ of the loop, namely, the Else clause. Now we return the global invariants:
$\left\{x y u_{0} z^{2}-u_{0}^{2} z^{2}+x y z u+x y u^{2}-x y z-2 x y u+x y, \ldots\right\}$ So, $x y u_{0} z^{2}-u_{0}^{2} z^{2}+x y z u+x y u^{2}-x y z-2 x y u+x y=0$ is one typical invariant that can be generated. Once again, here there are no need for Gröbner basis computation and the complexity of the described steps remains polynomial.

Example 9 illustrate our method for the case where the loop contains two conditional statements. In the presence of nested loops, our method generates ideals for invariants for each inner-loop and then generates a global invariant.

## 9. Experiments

The third column in Table 1a summarizes the type of linear algebraic problems associated with each kind of consecution approximation, listed in the second column, and with the semantic of the program instructions appearing in the first column. The last column in Table la gives some existential results which, we note, can also be used by other constraint-based approaches or reachability analysis methods. We have also used it to obtain some experimental results that attest to the effectiveness and scope of our methods considering the computation

Table 1. Examples and experimental results
(a) Linear algebraic problems and consecution approximations

| Prog. Loop | Aprox.Consec. | Linear Algebra | Existence Cond. |
| :--- | :--- | :--- | :--- |
| Affine/lin. inst. | Strong Scaling | Nullspaces | $\operatorname{Dim} \operatorname{Ker}\left(M_{D}\right) \geq 2$ <br> for any init. cond., <br> and $\operatorname{Ker}\left(M_{D}\right) \neq \emptyset$ <br> otherwise. |
| Affine/lin. inst. | Lambda Scaling | Eigenspaces | Dim_Eigen $\left(M_{D}\right) \geq 2$ <br> for any init. cond., <br> and Eigen $\left(M_{D}\right) \neq \emptyset$ <br> otherwise. |
| Algebraic/poly. inst. | Polynomial Scaling | Nullspaces | $\operatorname{Ker}\left(M_{D}-L_{T}\right) \geq 2$ <br> for any init. cond., <br> and $\operatorname{Ker}\left(M_{D}-L_{T}\right) \neq \emptyset$ <br> otherwise. |
|  |  |  | Dim_Ker $\left(M_{\Pi}-L_{T}\right) \geq 2$ <br> for any init. cond., <br> and $\operatorname{Ker}\left(M_{\square}-L_{T}\right) \neq \emptyset$ <br> otherwise. |
| Fractional inst. | Fractional Scaling | Nullspaces |  |
|  |  |  |  |

(b) Experimental results: computation of nullspaces and eigenspaces

| Loop prog. | Var. | Par. | Scaling | Basis inv. | CPU (s) |
| :---: | :---: | :---: | :---: | :---: | :---: |
| $1-\left[\begin{array}{c}x_{1}^{\prime}=2 x_{1}+x_{2}+1 \\ x_{2}^{\prime}=3 x_{2}+4\end{array}\right]$ | $\left\{x_{1}, x_{2}\right\}$ |  | $\begin{aligned} & \lambda \in\{9,6,4, \\ & 3,2,1\} \end{aligned}$ | $\left\{\left\{x_{2}^{2}+4 x_{2}+4\right\} ; \ldots\right\}$ | 0.39 |
| $2-\left[\begin{array}{l}{\left[s^{\prime}=s+i ; j^{\prime}=j-1 ;\right.} \\ \left.i^{\prime}=i ; j_{0}^{\prime}=j_{0}\right]\end{array}\right.$ | $\left\{s, i, j, j_{0}\right\}$ |  | $\lambda \in\{1\}$ | $\left\{\left\{j i+s, i^{2}, i j_{0}, i, j_{0}, 1\right\}\right\}$ | 1.27 |
| $3-\left[\begin{array}{l}x_{1}^{\prime}=a x_{1}+b x_{2}+c \\ x_{2}^{\prime}=d x_{1}+e x_{2}+f\end{array}\right]$ | $\left\{x_{1}, x_{2}\right\}$ | $\begin{aligned} & \{a, b, c, \\ & d, e, f\} \end{aligned}$ | $\begin{aligned} & \lambda \in\left\{0, d^{2}+\frac{1}{2} d e+\frac{1}{2} e^{2}\right. \\ & \pm \frac{1}{2} \operatorname{sqrt}\left(8 d^{3} e+d^{2} e^{2}\right. \\ & \left.\left.+6 d e^{3}+e^{4}\right), \ldots\right\} \\ & \hline \end{aligned}$ | $\left\{\left\{x_{1}-1\right\} ; \ldots\right\}$ | 1.35 |
| $4-\begin{aligned} & \left(\tau_{1}\right):\left[r^{\prime}=r+1 ; w^{\prime}=0 ;\right. \\ & \left.k^{\prime}=k-c_{1} ; c_{1}=c_{1}\right] \end{aligned}$ | $\left\{r, w, k, c_{1}\right\}$ |  | $\lambda \in\{0,1\}$ | $\begin{aligned} & \left\{\left\{r w, w^{2}, w k,\right.\right. \\ & \left.w c_{1}, w, k+c_{1}\right\} ; \\ & \left.\left\{c_{1}^{2}, c_{1}, 1\right\}\right\} \\ & \hline \end{aligned}$ | 0.37 |
| $\begin{aligned} & \left(\tau_{2}\right):\left[r^{\prime}=0 ; w^{\prime}=w+1 ;\right. \\ & k^{\prime}=k-c_{2} ; \\ & \left.c_{2}=c_{2}\right] \end{aligned}$ | $\left\{r, w, k, c_{2}\right\}$ |  | $\lambda \in\{0,1\}$ | $\begin{aligned} & \left\{\left\{r^{2}, r w, r k, r c_{2}, r\right\} ;\right. \\ & \left.\left\{w c_{2}+k\right\}\right\} \end{aligned}$ | 0.4 |
| $6-\begin{aligned} & \left(\tau_{3}\right):\left[r^{\prime}=r-1 ; w^{\prime}=0 ;\right. \\ & \left.k^{\prime}=k+c_{1} ; c_{1}=c_{1}\right] \end{aligned}$ | $\left\{r, w, k, c_{1}\right\}$ |  | $\lambda \in\{0,1\}$ | $\begin{aligned} & \left\{\left\{r w, w^{2}, w k, w c_{1}, w\right\} ;\right. \\ & \left.\left\{r c_{1}+k, c_{1}^{2}, c_{1}, 1\right\}\right\} . \end{aligned}$ | 0.39 |
| $7-\begin{aligned} & \left(\tau_{4}\right):\left[r^{\prime}=0 ; w^{\prime}=w-1 ;\right. \\ & \left.k^{\prime}=k+c_{2} ; c_{2}=c_{2}\right] \end{aligned}$ | $\left\{r, w, k, c_{2}\right\}$ |  | $\lambda \in\{0,1\}$ | $\begin{aligned} & \left\{\left\{r^{2}, r w, r k, r c_{2}, r\right\} ;\right. \\ & \left.\left\{c_{2}^{2}, w, c_{2}, 1\right\}\right\} . \end{aligned}$ | 0.43 |
| $\left.8-\left\lvert\, \begin{array}{c}x^{\prime}=x y+x \\ y^{\prime}=y^{2}\end{array}\right.\right]$ | $\{x, y\}$ |  | $T(x, y)=y^{2}+2 y+1$ | $\begin{aligned} & \left\{\left\{x^{2}, x y-x,\right.\right. \\ & \left.\left.y^{2}-2 y+1\right\}\right\} \\ & \hline \end{aligned}$ | 0.4 |
| $9-\left[\begin{array}{c}x^{\prime}=x y+x \\ y^{\prime}=y^{2}\end{array}\right]$ | $\{x, y\}$ |  | $T(x, y)=y^{2}+y$ | $\left\{\left\{x y, y^{2}-y\right\}\right\}$ | 0.41 |
| $10-\left[\begin{array}{c} x^{\prime}=x y+x \\ y^{\prime}=y^{2} \end{array}\right]$ | $\{x, y\}$ |  | $T(x, y)=y^{2}$ | $\left\{\left\{y^{2}\right\}\right\}$ | 0.41 |
| $11-\left[\begin{array}{c} x^{\prime}=x y+x \\ y^{\prime}=y^{2} \end{array}\right]$ | $\{x, y\}$ |  | $T(x, y)=y+1$ | \{ $\{x, y-1\}\}$ | 0.47 |
| $12-\left[\begin{array}{c} x^{\prime}=x y+x \\ y^{\prime}=y^{2} \end{array}\right]$ | $\{x, y\}$ |  | $T(x, y)=y$ | $\{\{y\}\}$ | 0.37 |
| 13- $\left.\left\lvert\, \begin{array}{c}x_{1}^{\prime}=\frac{x_{2}}{\left(x_{1}+x_{2}\right)} \\ x_{2}^{\prime}=\frac{x_{1}}{\left(x_{1}+2 x_{2}\right)}\end{array}\right.\right]$ | $\left\{x_{1}, x_{2}\right\}$ |  | $\begin{aligned} & T\left(x_{1}, x_{2}\right)=x_{1}^{2} \\ & +3 x_{1} x_{2}+2 x_{2}^{2} \end{aligned}$ | $\left\{\left\{x_{1} x_{2}\right\}\right\}$ | 0.37 |
| $14-\left[\begin{array}{l}x^{\prime}=a x^{2}+d x \\ y^{\prime}=a x y+d y\end{array}\right]$ | $\{x, y\}$ | $\{a, d\}$ | $T(x, y)=a^{2} x^{2}+2 a d x+d^{2}$ | $\left\{\left\{x^{2}, x y, y^{2}\right\}\right\}$ | 0.48 |
| 15- $\left.\begin{array}{c}x^{\prime}=a x^{2}+d x \\ y^{\prime}=a x y+i y^{2}+d y\end{array}\right]$ | $\{x, y\}$ | $\{a, d, i\}$ | $T(x, y)=a^{2} x^{2}+2 a d x+d^{2}$ | $\left\{\left\{x^{2}\right\}\right\}$ | 0.42 |
| 16- $\left[\begin{array}{l}x^{\prime}=a x^{2}+b x y+d x \\ y^{\prime}=a x y+b y^{2}+d y\end{array}\right]$ | $\{x, y\}$ | $\{a, b, d\}$ | $T(x, y)=a x+b y+d$ | $\{\{x, y\}\}$ | 0.43 |
| $17-\left[\begin{array}{c}x^{\prime}=b x y+c y^{2}+x \\ y^{\prime}=y\end{array}\right]$ | $\{x, y\}$ | $\{b, c, d\}$ | $T(x, y)=1$ | $\left\{\left\{y^{2}, y, 1\right\}\right\}$ | 0.45 |
| $\begin{gathered} {\left[x^{\prime}=b x y+c y^{2}\right.} \\ +d x+e y+f ; \\ 18-d x=b x^{2}+c x y \\ \left.y^{\prime}=d x+e y+f\right] \end{gathered}$ | $\{x, y\}$ | $\begin{aligned} & \{b, c, d, \\ & e, f\} \end{aligned}$ | $T(x, y)=-b x-c y$ | $\{\{x-y\}\}$ | 0.54 |
| $\left[x^{\prime}=a x^{2}+b x y ;\right.$ <br> 19- <br> $y^{\prime}=g x^{2}+h x y$ <br> $\left.+i y^{2}+k y\right]$ | $\{x, y\}$ | $\begin{aligned} & \{a, b, g, \\ & h, i, k\} \end{aligned}$ | $T(x, y)=a x+b y$ | $\{\{x\}\}$ | 0.47 |
| $20-\left[\begin{array}{l} x^{\prime}=a x^{2}+b x y \\ y^{\prime}=a x y+b y^{2} \end{array}\right]$ | $\{x, y\}$ | $\{a, b\}$ | $T(x, y)=a x+b y$ | $\{\{x, y\}\}$ | 0.46 |
| $21-\left[\begin{array}{l}x^{\prime}=a x^{2}+c y^{2}+f \\ y^{\prime}=a x^{2}+c y^{2}+l\end{array}\right]$ | $\{x, y\}$ | $\begin{aligned} & \{a, c, f, \\ & l\} \\ & \hline \end{aligned}$ | $T(x, y)=0$ | $\left\{\left\{x^{2}, 2 *((f * g \ldots\} ; \ldots\}\right.\right.$ | 1.84 |

of eigenspaces or nullspaces. By using efficient mathematical packages, e.g. Sage [SJ05], Maple, Mathematica, Lapack or Eispack, one can obtain the eigenvalues as closed-form algebraic expressions. That is, algebraic numbers which comprise the solutions of an algebraic equation in terms of its coefficients, relying only on $+,-, *, /$, and the extraction of roots. Also, eigenvalues are already obtained as algebraic numbers in practice, for large well known classes of matrices, such as when $n<5$, and when the matrix is $5 \times 5$ block triangular, among others. From Table 1a, we note that the computation of nullspaces or eigenspaces remain the main computational steps in our approaches. Table Sec-experimentb lists some of these experimental results focusing on the type of systems, scaling and basis of invariants that one could expect using our approaches. The first column refers to the specific problem of the experiment. The second column provides the numbers of variables. The third column gives the parameters used to represent a program class. The column Scaling shows the approximation of the consecution conditions. The column Basis Inv. presents the types of basis of invariants that one can generate. The last column refers to the cpu time required to compute those nullspaces or eigenspaces that turn out as vector spaces of invariants. We can see that our methods efficiently handle a large number of non-linear examples treated elsewhere in the literature. The experiment 2 is from [SSM04b] and the experiments $4-7$ relate to the loop transitions of the generalized readers-writers case studies from [SSM04b]. Experiments 8-20, listed in Table 1b, involve non-linear systems most of which can be shown to be beyond the limits of other recent approaches. The experiments $8-11$, expose the types of basis of invariants that one can generate considering four different polynomial scalings. More important experiment 3 and experiments $13-20$ refer to generic programs, i.e., large classes of programs and the associated generic basis of invariants that we provide. We use parameters to represent these classes of programs and we generate generic basis of invariants. Those parameters could also be use to abstract away some variables of a larger program. The experiment 12 involves a fractional system. We used the very complete Sage [SJ05] algebraic framework with interfaces written in Python so as to be able to access useful mathematical packages. Although the main contribution of this work is theoretical-we present theorems that could well be used with, or complement, other existing formal methods-the computation of these specific nullspaces or eigenspaces was conducted and depended on the Sage's on-line servers available. Those results show the strength of our approach for generating non-linear invariants for non-linear systems.

## 10. Discussion

The notions of Gröbner bases and their computations, together with the ideal membership problem are central to most recent approaches to program verification and static analysis [SSM04b, BBGL00, RCK07a, BLS96, CXYZ07, Kov08, KJ06, Cou05, MOS02, RCK07b, GT08, PC08]. In order to better understand the difficulties they incur, we first need some details on Gröbner basis and the ideal membership problem. Consider a multivariate polynomial, $Q=\sum_{i_{1}, \ldots, i_{n}} a_{i_{1}, \ldots, i_{n}} x_{1}{ }^{i_{1}} \ldots x_{n}{ }^{i_{n}}$, where the coefficients $a_{i_{1}, \ldots, i_{n}}$ are in a field $K$. How do we know if it is in an ideal $I$ of $K\left[X_{1}, \ldots, X_{n}\right]$ ? This is known as the Ideal membership problem. To handle it we can use a Gröbner basis $G=\left\{g_{1}, \ldots, g_{s}\right\}$ for $I$. There are algorithms that compute such bases as long as we know a finite generating basis for $I$ [Buc96, Fau99]. Then, we can compute the normal form of $Q$ for $I$ using the basis $G$. Denote the normal form by $N F_{G}(Q)$. We note that the use of a Gröbner basis guarantees the confluence and termination of those reductions. In general, we have $N F_{G}(Q)=\sum_{i_{1}, \ldots, i_{n}} f(a)_{i_{1}, \ldots, i_{n}} x_{1}{ }_{1} \ldots x_{n}{ }^{i_{n}}$, where $f(a)_{i_{1}, \ldots, i_{n}}$ is a combinations of the coefficients $a_{i_{1}, ., i_{n}}$. Then the statement $(Q \in I)$ is equivalent to the assertion $\left(N F_{G}(Q)=0\right)$, that is, all the coefficients $f(a)_{i_{1}, \ldots, i_{n}}$ are null.

Returning to the mentioned approaches for program verification and static analysis, the loop instructions are considered in order to form varieties and to build associated algebraic assertions and the ideal $I$. Then, these techniques compute a Gröbner basis $G$ for $I$. Next, they postulate a template polynomial $Q$, i.e., a polynomial with unknown coefficients, as a candidate invariant. As we have seen just above, $Q$ is an invariant if it belongs to the ideal $I$ or, in other words, if $\left(N F_{G}(Q)=0\right)$. So, the next step in these techniques, is to obtain the reduction $N F_{G}(Q)$. An important obstacle faced at this point is that all known algorithms for computing Gröbner basis and for the construction of the normal form reduction $N F_{G}(Q)$ are of doubly exponential time complexity. Having the normal form $N F_{G}(Q)$, they generate the set of candidate invariant constraints in the form of the system of equations $\left(N F_{G}(Q)=0\right)$, and attempt to solve it directly. Moreover, as we have seen in Sect. 4.3, as soon as the loop contains a non-linear instruction, the candidate invariant constraints results in a non-linear system of equation, which makes its resolution all but unfeasible. Further, there are no conditions over the degree of their candidate invariants that would guarantee the non-triviality of the resulting invariant, when it can be computed.

In terms of performance and efficiency, we succeeded in reducing the non-linear loop invariant generation problem to a linear algebraic problem, i.e., to the computation of eigenspaces of specific morphisms. Our techniques have fewer computational steps: we first compute some specific matrices and then we compute their nullspaces. Each of these steps remains of polynomial time complexity. Further, our approaches do not simply generate an invariant at a time. Instead, we generate an ideal of invariants which is a large-infinite-structure. We also handle fractional systems and our algorithm incorporates a strategy to find degree bounds which allow for the automatic generation of ideals of non-trivial invariants. Moreover, as one of the main results, we provide very general sufficient conditions allowing for the existence and computation of such invariant ideals. Note that these conditions could be directly used by any other invariant generation methods.

As a more applied motivation, our techniques can be made to bear on new domains that require the computation of complex invariants. Along these lines, some recent work on security [RM11c, RM09, RM11b, RM11a], showed how such invariants play a central role in static analysis of malwares, e.g., viruses, and how they can be used to build new invariant-based intrusion detection systems. Invariants generated over malware codes are strong semantic aware signatures that can be used to analyze and identify intrusions caused by such malicious code. These new approaches could form parts of intrusion detection systems where an alarm is a proof of abnormal behavior caused by the violation of a precomputed invariant induced by the intrusion. We note that binary code gives rise to non-linear arithmetic and the methods described here allow, as we have shown, for the generations of complex and precise invariants in such cases. We stress that the more the complex the invariant is, the harder it will be to morph the corresponding signatures in an automatic way.

## 11. Conclusions

Our primary goal and motivation were to provide invariant generation methods for static analysis that could serve as a basis for automatic program verification.

We have shown that the preconditions for discrete transitions can be viewed as morphisms over a vector space of bounded degree polynomials. These morphisms, in turn, could be suitably represented by matrices. By doing so, we succeeded in reducing the non-linear loop invariant generation problem to linear algebraic problems or, more precisely, to the computation of eigenspaces of these morphisms. We also treated fractional systems and our algorithms incorporate a strategy to find degree bounds for candidate invariants, thus allowing for the automatic generation of non-trivial invariants.

These techniques lead to algorithms of much lower time complexity than other modern approaches. The latter incur in computations which are of a doubly exponential time complexity while, by contrast, our techniques induce algorithms of polynomial time complexity.

Further, our methods do not generate a single invariant at a time. Instead, we generate non-linear invariant ideals, which are infinite structures, giving rise to families of non-trivial invariants. As another important main result, we provided very general sufficient conditions that can guarantee the existence of such invariant ideals.

We also noted that our techniques could be combined with other formal verification methods and their associated tools. A case in point are formal methods that treat logics with uninterpreted functions [GT06], which can handle function calls and operating system calls.

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## A. Appendix

Proof of Theorem 1 If $Q\left(X_{1}^{\prime}, \ldots, X_{n}^{\prime}\right)-\lambda Q\left(X_{1}, \ldots, X_{n}\right)$ belongs to the ideal $I$ generated by the family ( $X_{1}^{\prime}-$ $\left.L_{1}, \ldots, X_{n}^{\prime}-L_{n}\right)$, then there exists a family $\left(A_{1}, \ldots, A_{n}\right)$ of polynomials in $\mathbb{R}\left[X_{1}^{\prime}, \ldots, X_{n}^{\prime}, X_{1}, \ldots, X_{n}\right]$ such that $Q\left(X_{1}^{\prime}, \ldots, X_{n}^{\prime}\right)-\lambda Q\left(X_{1}, \ldots, X_{n}\right)=\left(X_{1}^{\prime}-L_{1}\right) A_{1}+\ldots+\left(X_{n}^{\prime}-L_{n}\right) A_{n}$. Letting $X_{i}^{\prime}=L_{i}$, we obtain $Q\left(L_{1}\left(X_{1}, \ldots, X_{n}\right), \ldots\right.$,
$\left.L_{n}\left(X_{1}, \ldots, X_{n}\right)\right)=\lambda Q\left(X_{1}, \ldots, X_{n}\right)$. Conversely suppose
$Q\left(L_{1}\left(X_{1}, . ., X_{n}\right), . ., L_{n}\left(X_{1}, . ., X_{n}\right)\right)=\lambda Q\left(X_{1}, . ., X_{n}\right)$. Then as $Q\left(X_{1}^{\prime}, \ldots, X_{n}^{\prime}\right)$ is equal to $Q\left(L_{1}, . ., L_{n}\right)$ modulo the ideal $I$, we get $Q\left(X_{1}^{\prime}, . ., X_{n}^{\prime}\right)=\lambda Q\left(X_{1}, \ldots, X_{n}\right)$ modulo $I$.

Proof of Theorem 2 Let $Q$ be a polynomial in $\mathbb{R}_{r}\left[X_{1}, \ldots, X_{n}\right]$. We have the following sequence of deduction: $\left(Q\left(L_{1}\left(X_{1}, \ldots, X_{n}\right), . .\right.\right.$,
$\left.\left.L_{n}\left(X_{1}, \ldots, X_{n}\right)\right)=\lambda Q\left(X_{1}, \ldots, X_{n}\right)\right) \Leftrightarrow(\mathscr{M}(Q)=\lambda Q) \Leftrightarrow(\mathscr{M}(Q)=\lambda \operatorname{Id}(Q)) \Leftrightarrow((\mathscr{M}-\lambda \operatorname{Id})(Q)$ $\left.=0_{\mathbb{R}}\left[X_{1}, \ldots, X_{n}\right]\right) \Leftrightarrow(Q \in \operatorname{Ker}(M-\lambda I))$. Using the definition of an invariant and Theorem 1 , we can see that $Q$ will be a $\lambda$-scale invariant if and only if it belongs to the eigenspace corresponding to $\lambda$.

Proof of Corollary 1 Suppose $M$ is block triangular with blocks $4 \times 4$ or less, then its characteristic polynomial will a product of polynomials of degree less than four, whose roots can be calculated by the Lagrange resolvent method [Lan02]. For the second assertion, we already know that 1 is an eigenvalue. Suppose that the corresponding eigenspace is of dimension exactly one. Then the only vectors in that space are the constant polynomials. If it is of dimension two or more, then we get non-trivial polynomials in the eigenspace.

Proof of Theorem 3 We first consider Theorem 2. The initiation step defines on $\mathbb{R}_{r}\left[x_{1}, \ldots, x_{n}\right]$ a linear form on this space, namely, $I_{u}: P \mapsto P\left(u_{1}, \ldots, u_{n}\right)$. Hence, initial values correspond to a hyperplane of $\mathbb{R}_{r}\left[X_{1}, . ., X_{n}\right]$ given by the kernel of $I_{u}$, which is $\left\{Q \in \mathbb{R}_{r}\left[X_{1}, . ., X_{n}\right] \mid Q\left(u_{1}, \ldots, u_{n}\right)=0\right\}$. If we add initial conditions of the form $\left(x_{1}(0)\right.$ $\left.=u_{1}, \ldots, x_{n}(0)=u_{n}\right)$, we are looking for a $\lambda$-scale invariant in $\mathbb{R}_{r}\left[x_{1}, \ldots, x_{n}\right]$ that belongs to the hyperplane $P\left(u_{1}, \ldots, u_{n}\right)=0$, i.e., we are looking for $Q$ in $\operatorname{ker}(M-\lambda I) \cap\left\{P \mid P\left(u_{1}, \ldots, u_{n}\right)=0\right\}$.

Proof of Corollary 2 We take each direction, in turn. $[(\Rightarrow)]$ There is a $\lambda$-scale invariant for any initial value. Then the corresponding eigenspace has dimension at least 2 . Indeed, if the space was of dimension only 1 (which is at least necessary to have $\lambda$-invariants), than taking any nonzero vector $Q$ in the eigenspace (i.e. a $\lambda$-invariant), $Q$ should lie in any hyperplane of initial values. That is for every $n$-tuple $\left(u_{1}, \ldots, u_{n}\right)$ one would have $Q\left(u_{1}, \ldots, u_{n}\right)=0$, hence $Q=0$, which is absurd. $[(\Leftarrow)]$ Any eigenspace of $M$ with dimension at least 2 will intersect any space, in particular any hyperplane, given by any initial constraints. As any hyperplane is of co-dimension one in $\mathbb{R}_{r}\left[X_{1}, \ldots, X_{n}\right]$, it must have a nonzero intersection with any subspace of dimension strictly greater than one. This establishes the result.

Proof of Theorem $4[(\Rightarrow)]$ If $Q\left(X_{1}^{\prime}, . ., X_{n}^{\prime}\right)-T Q\left(X_{1}, . ., X_{n}\right)$ belongs to the ideal $I$ generated by the family $\left(X_{1}^{\prime}-P_{1}, \ldots, X_{n}^{\prime}-P_{n}\right)$, then there exists a family $\left(A_{1}, \ldots, A_{n}\right)$ of polynomials in $\mathbb{R}\left[X_{1}^{\prime}, . ., X_{n}^{\prime}, X_{1}, . ., X_{n}\right]$ such that $Q\left(X_{1}^{\prime}, . ., X_{n}^{\prime}\right)$
$-\lambda Q\left(X_{1}, \ldots, X_{n}\right)=\left(X_{1}^{\prime}-P_{1}\right) A_{1}+\ldots+\left(X_{n}^{\prime}-P_{n}\right) A_{n}$. Letting $X_{i}^{\prime}=P_{i}$, we obtain $Q\left(P_{1}\left(X_{1}, \ldots, X_{n}\right), \ldots\right.$, $\left.P_{n}\left(X_{1}, \ldots, X_{n}\right)\right)=T Q\left(X_{1}, \ldots, X_{n}\right)$. $[(\Leftarrow)]$ Conversely suppose
$Q\left(P_{1}\left(X_{1}, . ., X_{n}\right), . ., P_{n}\left(X_{1}, . ., X_{n}\right)\right)=T Q\left(X_{1}, . ., X_{n}\right)$. Then as $Q\left(X_{1}^{\prime}, . ., X_{n}^{\prime}\right)$ is equal to $Q\left(P_{1}, . ., P_{n}\right)$ modulo the ideal $I$, we get $Q\left(X_{1}^{\prime}, \ldots, X_{n}^{\prime}\right)=\lambda Q\left(X_{1}, \ldots, X_{n}\right)$ modulo $I$. This establishes the result.

Proof of Theorem 5 Let $Q$ be a polynomial in $\mathbb{R}\left[X_{1}, \ldots, X_{n}\right]$. Such a polynomial $Q$ is $T$-invariant if and only if $Q\left(P_{1}\left(X_{1}, . ., X_{n}\right), . .\right.$,
$\left.P_{n}\left(X_{1}, \ldots, X_{n}\right)\right)=T\left(X_{1}, \ldots, X_{n}\right) Q\left(X_{1}, \ldots, X_{n}\right)$, i.e., if and only if $\mathscr{M}(Q)=\mathscr{L}(Q) \Leftrightarrow(\mathscr{M}-\mathscr{L})(Q)=0_{\mathbb{R}\left[X_{1}, \ldots, X_{n}\right]}$. Writing this in matrix equivalent terms we have $((M-L) Q=0) \Leftrightarrow(Q \in \operatorname{Ker}(M-L))$, and the result follows. $\square$

Proof of Theorem 6 From linear algebra, we know that $M-L$ with a non-trivial kernel is equivalent to it having rank strictly less than the dimension $v(r)$ of $\mathbb{R}_{r}\left[x_{1}, \ldots, x_{n}\right]$. This is equivalent to the fact that each $v(r) \times v(r)$ sub-determinant of $M_{D}-L_{T}$ is equal to zero. Those determinants are polynomials with variables $\left(t_{1}, \ldots, t_{v(d-1)}\right)$, which we will denote by $V_{1}\left(t_{1}, \ldots, t_{v(d-1)}\right), \ldots, V_{s}\left(t_{1}, \ldots, t_{v(d-1)}\right)$. From the form of $L$, this is zero when $\left(t_{1}, \ldots, t_{v(d-1)}\right)=(0, \ldots, 0)$. Hence, $M-L$ has its last column equal to zero, giving a common root for these polynomials, corresponding to the constant invariants.

Proof of Theorem 7 Consider Theorem 5. The initiation step defines on $\mathbb{R}_{r}\left[x_{1}, \ldots, x_{n}\right]$ a linear form on this space, namely, $I_{u}: P \mapsto P\left(u_{1}, \ldots, u_{n}\right)$. Thus, initial values correspond to a hyperplane of $\mathbb{R}_{r}\left[X_{1}, . ., X_{n}\right]$ given by the kernel of $I_{u}$, which is $\left\{Q \in \mathbb{R}_{r}\left[X_{1}, . ., X_{n}\right] \mid Q\left(u_{1}, \ldots, u_{n}\right)=0\right\}$. With initial conditions $\left(x_{1}(0)=u_{1}, \ldots, x_{n}(0)=u_{n}\right)$, we are looking for a $T$-scale differential invariant in $\mathbb{R}_{r}\left[x_{1}, \ldots, x_{n}\right]$ that belongs to the hyperplane $P\left(u_{1}, \ldots, u_{n}\right)=$ 0 , i.e., we are looking for $Q$ in $\operatorname{Ker}(M-L) \cap\left\{P \mid P\left(u_{1}, \ldots, u_{n}\right)=0\right\}$.

Proof of Corollary $3[(\Rightarrow)$ ] If there is a $T$-scale invariant for any initial value, then the corresponding eigenspace has dimension at least 2 . Indeed, if the space was of dimension only 1 (which is at least necessary to have $T$ invariants), taking any non-zero vector $Q$ in the eigenspace (i.e. a $T$-invariant), $Q$ should lie in any hyperplane of initial values, and so for every n-tuple $\left(u_{1}, \ldots, u_{n}\right)$, one would have $Q\left(u_{1}, \ldots, u_{n}\right)=0$, hence $Q=0$, which is
absurd. $[(\Leftarrow)]$ Any eigenspace of $M_{D}-L_{T}$ with dimension at least 2 will intersect any space given by any initial constraints. This establishes the result.
Proof of Corollary 4 The right singular vectors corresponding to vanishing singular values of $\overline{M-L}$ span the null space of $\bar{M}-L$. The left singular vectors corresponding to the non-zero singular values of $\overline{M-L}$ span the range of $\overline{M-L}$. As a consequence, the rank of $\overline{M-L}$ equals the number of non-zero singular values which is the same as the number of non-zero elements in the matrix $S$.
Proof of Theorem $8[(\Rightarrow)]$ If $Q\left(X_{1}^{\prime}, \ldots, X_{n}^{\prime}\right)-F Q\left(X_{1}, \ldots, X_{n}\right)$ belongs to the fractional ideal $J$ generated by the family $\left(X_{1}^{\prime}-P_{1} / Q_{1}, \ldots, X_{n}^{\prime}-P_{n} / Q_{n}\right)$, then there exists a family $\left(A_{1}, \ldots, A_{n}\right)$ of fractional functions in $\mathbb{R}\left(X_{1}^{\prime}, \ldots, X_{n}^{\prime}, X_{1}\right.$, .., $X_{n}$ ) such that $Q\left(X_{1}^{\prime}, \ldots, X_{n}^{\prime}\right)-F Q\left(X_{1}, \ldots, X_{n}\right)=\left(X_{1}^{\prime}-P_{1} / Q_{1}\right) A_{1}+\ldots+\left(X_{n}^{\prime}-P_{n} / Q_{n}\right) A_{n}$ Letting $X_{i}^{\prime}=\frac{P_{i}}{Q_{i}}$ we obtain $Q\left(\frac{P_{1}}{Q_{1}}, \ldots, \frac{P_{n}}{Q_{n}}\right)=\lambda Q\left(X_{1}, \ldots, X_{n}\right) .[(\Leftarrow)]$ Conversely suppose $Q\left(\frac{P_{1}}{Q_{1}}, \ldots, \frac{P_{n}}{Q_{n}}\right)=F Q\left(X_{1}, \ldots, X_{n}\right)$. Then as $Q\left(X_{1}^{\prime}, \ldots, X_{n}^{\prime}\right)$ is equal to $Q\left(\frac{P_{1}}{Q_{1}}, \ldots, \frac{P_{n}}{Q_{n}}\right)$ modulo the ideal $J$, we get that $Q\left(X_{1}^{\prime}, \ldots, X_{n}^{\prime}\right)=F Q\left(X_{1}, \ldots, X_{n}\right)$ modulo $J$. And we have the result.
Proof of Theorem 9 Let $Q$ be a polynomial in $\mathbb{R}\left[X_{1}, \ldots, X_{n}\right]$. In fact, a polynomial $Q$ is $T / \Pi^{r}$-invariant if and only if $Q\left(P_{1}\left(X_{1}, \ldots, X_{n}\right), \ldots, P_{n}\left(X_{1}, . ., X_{n}\right)\right)=T / \Pi^{r}\left(X_{1}, . ., X_{n}\right) Q\left(X_{1}, . ., X_{n}\right)$, which is equivalent to $\Pi^{r} Q\left(P_{1}\left(X_{1}, \ldots, X_{n}\right)\right.$,
$\left.\ldots, P_{n}\left(X_{1}, . ., X_{n}\right)\right)=T\left(X_{1}, \ldots, X_{n}\right) Q\left(X_{1}, \ldots, X_{n}\right)$, and this holds if and only if $(\mathscr{M}(Q)=\mathscr{L}(Q)) \Leftrightarrow$ $\left((\mathscr{M}-\mathscr{L})(Q)=0_{\mathbb{R}\left[X_{1}, \ldots, X_{n}\right]}\right.$ Writing this in matrix equivalent terms we get $((M-L) Q=0) \Leftrightarrow(Q \in \operatorname{Ker}(M-L))$, and the result follows.
Proof of Theorem 10 We first consider Theorem 9. The initiation step defines on $\mathbb{R}_{r}\left[x_{1}, \ldots, x_{n}\right]$ a linear form on this space, namely, $I_{u}: P \mapsto P\left(u_{1}, \ldots, u_{n}\right)$. Hence, initial values correspond to a hyperplane of $\mathbb{R}_{r}\left[X_{1}, . ., X_{n}\right]$ given by the kernel of $I_{u}$, which is $\left\{Q \in \mathbb{R}_{r}\left[X_{1}, \ldots, X_{n}\right] \mid Q\left(u_{1}, \ldots, u_{n}\right)=0\right\}$. With initial conditions $\left(x_{1}(0)=\right.$ $\left.u_{1}, \ldots, x_{n}(0)=u_{n}\right)$, we are looking for a strong-scale differential invariant in $\mathbb{R}_{r}\left[x_{1}, \ldots, x_{n}\right]$ that belongs to the hyperplane $P\left(u_{1}, \ldots, u_{n}\right)=0$, i.e., we are looking for $Q$ in $\operatorname{Ker}(M-L) \cap\left\{P \mid P\left(u_{1}, \ldots, u_{n}\right)=0\right\}$.
Proof of Corollary $5[(\Rightarrow)]$ If there is a non-trivial $F$-scale invariant for any initial value, then the corresponding eigenspace has dimension at least 2 . Indeed, if the space was of dimension only 1 (which is at least necessary to have $F$-invariants), taking any non-zero vector $Q$ in the eigenspace (i.e. a $F$-invariant), $Q$ should lie in any hyperplane of initial values, i.e. for every n-tuple ( $u_{1}, \ldots, u_{n}$ ), one would have $Q\left(u_{1}, \ldots, u_{n}\right)=0$, hence $Q=0$, which is absurd. $[(\Leftarrow)]$ Any intersection between an eigenspace of $M$ with dimension at least 2 will intersect any space given by any initial constraints. And we have the result.
Proof of Theorem 11 Let $f_{1}^{(j)}, \ldots, f_{n_{j}}^{(j)} \in K\left[X_{1}, \ldots, X_{n}\right]$ be such that $I_{j}=\left(f_{1}^{(j)}, \ldots, f_{n_{j}}^{(j)}\right)$, for all $j$ in $[1, k]$. Let $\beta \in$ $\left(\otimes\left(I_{1}, \ldots, I_{k}\right)\right)$ Then there exists $e_{1}, \ldots, e_{n_{1} n_{2} \ldots n_{k}} \in K\left[X_{1}, \ldots, X_{n}\right]$ such that $\beta=e_{1} \delta_{1}+\ldots+e_{n_{1} n_{2} \ldots n_{k}} \delta_{n_{1} n_{2} \ldots n_{k}}$. Also, by construction of $\otimes\left(I_{1}, \ldots, I_{k}\right)$ we know that for all $r \in\left[1, \ldots, n_{1} n_{2} \ldots n_{k}\right], \delta_{r} \in \otimes\left(I_{1}, \ldots, I_{k}\right)$. In other words, there is $\left(\alpha_{1}^{(r)}, \ldots, \alpha_{k}^{(r)}\right) \in I_{1} \times \cdots \times I_{k}$ such that $\delta_{r}=\Pi_{i=0}^{k} \alpha_{i}^{(r)}$. Then we have $\beta=\sum_{j=1}^{n_{1} n_{2} \ldots n_{k}}\left[\lambda_{j} \Pi_{i=1}^{k} \alpha_{i}^{(j)}\right]$. Now, for all $m$ in $[1, k]$, if $I_{m}$ correspond to a precomputed inductive ideal of invariants associated to a transition $\tau_{m}$ at location $l$, then for all $j \in\left[1, n_{1} n_{2} \ldots n_{k}\right]$ we have $\alpha_{m}^{(j)}(X 1, \ldots, X n)=0$. Hence, for all $j \in\left[1, n_{1} n_{2} \ldots n_{k}\right]$ we get $\Pi_{i=1}^{k} \alpha_{i}^{(j)}=0$. Finally, we obtain $\beta\left(X_{1}, \ldots, X_{n}\right)=0$ for all $m$ in $\left[1, n_{1} n_{2} \ldots n_{k}\right]$. In other words, $\beta\left(X_{1}, \ldots, X_{n}\right)=0$ is an algebraic assertion true at any step of the iteration of the loop for any transition $\tau_{m}$ that could possibly be taken. Then $\left(\beta\left(X_{1}, \ldots, X_{n}\right)=0\right)$ is an inductive invariant and we can conclude that $\left(\otimes\left(I_{1}, \ldots, I_{k}\right)\right)$ is an ideal of inductive invariants.

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