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# Exact linear reductions of dynamical models 

Alexander Demin*, Elizaveta Demitraki ${ }^{\dagger}$, and Gleb Pogudin ${ }^{\ddagger}$


#### Abstract

Dynamical models described by ordinary differential equations (ODEs) are a fundamental tool in the sciences and engineering. Exact reduction aims at producing a lower-dimensional model in which each macro-variable can be directly related to the original variables, and it is thus a natural step towards the model's formal analysis and mechanistic understanding. We present an algorithm which, given a polynomial ODE model, computes a longest possible chain of exact linear reductions of the model such that each reduction refines the previous one, thus giving a user control of the level of detail preserved by the reduction. This significantly generalizes over the existing approaches which compute only the reduction of the lowest dimension subject to an approach-specific constraint. The algorithm reduces finding exact linear reductions to a question about representations of finitedimensional algebras. We provide an implementation of the algorithm, demonstrate its performance on a set of benchmarks, and illustrate the applicability via case studies. Our implementation is freely available at https://github.com/x3042/ExactODEReduction.jl.


Key words. ordinary differential equations, exact reduction, lumping, dimensionality reduction, matrix algebras

MSC codes. 34C20, 34-04, 16G10

1. Introduction. Ordinary Differential Equations (ODEs) provide a powerful and expressive language for describing systems evolving in real-time and, thus, are widely used both in the sciences and engineering. This motivates development of formal methods to analyse the structure of models defined using ODEs. One important problem which has been studied actively in the past decade from this angle is model reduction [16, 17, 6, 7, 28].

In general, model reduction refers to a variety of techniques aiming at replacing the model of interest with a simpler one while preserving, at least approximately, some of the important features of the original model. Traditionally, approximate methods such as, e.g., balanced truncation [1] have been employed. Exact model reduction is a complementary approach in which one lowers the dimension of the model without introducing approximation errors. Such reductions are of particular interest in the context of performing formal analysis or deriving mechanistic insights.

In this paper, we will focus on an important class of such reductions, exact linear lumpings, which correspond to finding a self-consistent system of differential equations for a set of macrovariables in which each macro-variable is a linear combination of the original variables. The case of the macro-variables being sums of the original variables has been extensively studied, see e.g. [16, 8, 17, 7]. In particular, ERODE software has been developed [6] which efficiently finds the optimal partition of the original variables into macro-variables. A recent software CLUE [28] was a step towards lifting these restrictions on the macro-variables. Unlike the prior

[^0]approaches, the macro-variables produced by CLUE are allowed to involve any coefficients (not just zeroes and ones as before), and, thus, the dimension of the reduced model could be significantly lower [28, Table 1]. However, the input of the algorithm consisted not only of a model but also of linear forms in the state variables to be preserved (the observables). Such a set of observables may or may not be available, and guessing it correctly is crucial for finding low-dimensional reductions. In this paper, we aim at taking the best from both worlds: requiring only a model as input (as ERODE) and allowing the macro-variables to be any linear combinations of the original variables (as CLUE).

The main result of the paper is an algorithm for finding arbitrary exact linear reductions when given only a polynomial ODE model with rational coefficients. Note that the question of finding such an arbitrary linear lumping of the lowest possible dimension may not be the most meaningful one since any linear first integral yields a reduction of dimension one with constant dynamics. Instead, our algorithm finds a longest possible chain of lumpings in which each reduction refines the next one (for details, see Section 2) so that a user can choose the desired level of details to be preserved by reduction by moving along the chain and may find reductions which would likely be missed by ERODE and CLUE (e.g., see Example 2.4 and subsection 5.2 ). Such generality comes with a price: our software is typically slower than CLUE and ERODE.

Our algorithm is based on combining the connection of the linear lumping problem to the problem of finding a common invariant subspace of a set of matrices [25, 28] with the structure theory of finite-dimensional algebras. We use the general framework of existing algorithms over finite [33] and algebraically closed [9, 36] fields and achieve desired efficiency by

- sparsity-aware algorithm for finding a basis of an algebra (Subsection 3.2);
- exploiting the structure of the input to compute mostly with rational numbers and postponing passing to algebraic number fields as much as possible;
- using sparse linear algebra and modular computation to avoid large matrices and expression swell, respectively.
We implemented our algorithm, and the implementation is publicly available at https: //github.com/x3042/ExactODEReduction.jl. We evaluate its performance on a set of benchmarks from the BioModels database [27], a large collection of models from life sciences, and demonstrate the produced reduction for two case studies.

The rest of the paper is organized as follows. In Section 2, we give precise definition of exact linear reduction, and formulate explicitly the algorithmic problem we solve in the paper. Section 3 contains detailed description of the algorithm and its justification. We describe our implementation and report its performance in Section 4. Section 5 contains the case studies describing the reductions produced by our software.
2. Problem statement. In the paper, the transpose of a matrix $M$ is denoted by $M^{T}$. For a vector $\mathbf{x}=\left(x_{1}, \ldots, x_{n}\right)$ of indeterminates, by $\mathbb{C}[\mathbf{x}]$ (resp., $\left.\mathbb{R}[\mathbf{x}], \mathbb{Q}[\mathbf{x}]\right)$ we will denote the set of polynomials in $\mathbf{x}$ with complex (resp., real, rational) coefficients.

Definition 2.1 (Lumping). Consider a system of ODEs with polynomial right-hand side of the form

$$
\begin{equation*}
\mathbf{x}^{\prime}=\mathbf{f}(\mathbf{x}) \tag{2.1}
\end{equation*}
$$

where $\mathbf{x}=\left(x_{1}, \ldots, x_{n}\right), \mathbf{f}=\left(f_{1}, \ldots, f_{n}\right)$, and $f_{1}, \ldots, f_{n} \in \mathbb{C}[\mathbf{x}]$. We say that a linear transformation $\mathbf{y}=\mathbf{x} L$ with $\mathbf{y}=\left(y_{1}, \ldots, y_{m}\right), L \in \mathbb{C}^{n \times m}$, and $\operatorname{rank} L=m$ is a lumping of (2.1) if there exist $\mathbf{g}=\left(g_{1}, \ldots, g_{m}\right)$ with $g_{1}, \ldots, g_{m} \in \mathbb{C}[\mathbf{y}]$ such that

$$
\mathbf{y}^{\prime}=\mathbf{g}(\mathbf{y})
$$

for every solution $\mathbf{x}$ of (2.1). We say that $m$ is the dimension of the lumping. The variables $\mathbf{y}$ in the reduced system are called macro-variables.

Throughout this section, we will work with the following running example [30, Example 1]. We will consider a chemical reaction network consisting of

- A chemical species $X$.
- Species $A_{U U}, A_{U X}, A_{X U}$, and $A_{X X}$. Each of them is one of the states of a molecule $A$ with two identical binding sites, which can be either unbound ( U in the subscript) or bound (X in the subscript) to $X$.
For simplicity, we will assume that all the reaction rates are equal to one. The dynamics of the network is defined by the following reactions:

$$
\begin{equation*}
X+A_{U *} \rightleftarrows A_{X *} \quad X+A_{* U} \rightleftarrows A_{* X} \tag{2.2}
\end{equation*}
$$

where $* \in\{X, U\}$. Under the laws of the mass-action kinetics [14, Ch. 7], the reactions (2.2) yield the following ODE system (where $[S]$ denotes the concentration of the species $S$ ):

$$
\left\{\begin{array}{l}
{[X]^{\prime}=\left[A_{X U}\right]+\left[A_{U X}\right]+2\left[A_{X X}\right]-[X]\left(\left[A_{X U}\right]+\left[A_{U X}\right]+2\left[A_{U U}\right]\right)}  \tag{2.3}\\
{\left[A_{U U}\right]^{\prime}=\left[A_{X U}\right]+\left[A_{U X}\right]-2[X]\left[A_{U U}\right]} \\
{\left[A_{U X}\right]^{\prime}=\left[A_{X X}\right]+[X]\left[A_{U U}\right]-[X]\left[A_{U X}\right]-\left[A_{U X}\right]} \\
{\left[A_{X U}\right]^{\prime}=\left[A_{X X}\right]+[X]\left[A_{U U}\right]-[X]\left[A_{X U}\right]-\left[A_{X U}\right]} \\
{\left[A_{X X}\right]^{\prime}=[X]\left[A_{X U}\right]+[X]\left[A_{U X}\right]-2\left[A_{X X}\right]}
\end{array}\right.
$$

Example 2.2 (Conservation laws as lumpings). We show that matrix $L=\left(\begin{array}{lllll}0 & 1 & 1 & 1 & 1\end{array}\right)^{T}$ yields a lumping of (2.3). We have

$$
y=\left(\left[\begin{array}{lll}
{[X]} & {\left[A_{U U}\right]} & {\left[A_{U X}\right]}
\end{array}\left[A_{X U}\right] \quad\left[A_{X X}\right]\right) \cdot L=\left[A_{U U}\right]+\left[A_{U X}\right]+\left[A_{X U}\right]+\left[A_{X X}\right]\right.
$$

Using (2.3) one can check that $y^{\prime}=0$, so we can take $g(y)=0$. Indeed, $y$ is the total concentration of $A$ and must be constant. Furthermore, any linear conservation law yields a lumping of dimension one.

Example 2.3 (More informative lumping). Another lumping for the same system (2.3) is given by the matrix

$$
L=\left(\begin{array}{lll}
1 & 0 & 0 \\
0 & 2 & 0 \\
0 & 1 & 1 \\
0 & 1 & 1 \\
0 & 0 & 2
\end{array}\right) \Longrightarrow\left\{\begin{array}{l}
y_{1}=[X] \\
y_{2}=2\left[A_{U U}\right]+\left[A_{U X}\right]+\left[A_{X U}\right] \\
y_{3}=\left[A_{U X}\right]+\left[A_{X U}\right]+2\left[A_{X X}\right]
\end{array}\right.
$$

The macro-variables will satisfy a self-contained system

$$
y_{1}^{\prime}=y_{3}-y_{1} y_{2}, \quad y_{2}^{\prime}=y_{3}-y_{1} y_{2}, \quad y_{3}^{\prime}=-y_{3}+y_{1} y_{2}
$$

The rationale behind this reduction is that $y_{2}$ and $y_{3}$ are "concentrations" of unbound and binded sites, respectively.

The above examples demonstrate that one system can have several lumpings (in fact, (2.3) has more), so a natural question is how to find useful lumpings. The state-of-the-art software tools CLUE [28] and ERODE [6] approach this question by finding the lumping of the smallest dimension satisfying certain constraints:

- preserving some quantities of interest unlumped (for CLUE [28]);
- or coming from a partition of the state variables (for ERODE [6]).

Both constraints may be too restrictive: not all interesting lumpings come from a partition (like the ones in Examples 2.2 and 2.3; see also [28, Table 1]), and it may be complicated to guess in advance meaningful quantities to preserve.

Example 2.4 (Example hard for CLUE and ERODE, see also Subsection 5.2). Consider another chemical reaction network [15, Eq. (19.20)] (originally due to Daniel Knight):

$$
\begin{array}{cc}
E+S \rightleftarrows E S & \\
\downarrow & \\
& E \neq P \\
E^{*}+S
\end{array} \quad \begin{aligned}
& \\
&
\end{aligned}
$$

As in the case of (2.2), we transform the reactions into an ODE system using the law of mass-action kinetics taking all the rates to be one ${ }^{1}$ :

$$
\left\{\begin{array}{l}
{[E]^{\prime}=2[E S]+\left[E^{*}\right]-[E][S]-[E][P]-[E],} \\
{[S]^{\prime}=2[E S]-[E][S],} \\
{[P]^{\prime}=[E S]-[E][P],} \\
{[E S]^{\prime}=[E][S]+[E][P]-3[E S],} \\
{\left[E^{*}\right]^{\prime}=[E]+[E S]-\left[E^{*}\right]}
\end{array}\right.
$$

One meaningful linear reduction is $y=[E]+[E S]-\left[E^{*}\right]$ with the equation $y^{\prime}=-2 y$. The macro-variable $y$ can be understood as a potential between the amount of $E$ (typically enzyme), both in the free form $E$ and as a part of the complex $E S$, and $E^{*}$ (typically inactivated enzyme). This reduction does not come from a subdivision of the species, so it cannot be found by ERODE. Furthermore, finding it using CLUE would require knowing this macro-variable in advance.

An alternative approach would be to find all the lumpings and let the user to choose which ones to use. The problem is that there may be easily an inifinite number of lumpings, for example, similarly to Example 2.2, one can show that the matrix

$$
L=\left(\begin{array}{lllll}
\alpha & 1 & 1+\alpha & 1+\alpha & 1+2 \alpha
\end{array}\right)^{T}
$$

[^1]yields a lumping of (2.3) for every number $\alpha$. Furthermore, as we will explain later, the lumpings are in a bijection with the invariant subspace of certain matrices coming from the Jacobian of $\mathbf{f}(\mathbf{x})$ and (at least for arbitrary matrices) the invariant subspaces can form an arbitrary algebraic variety [32].

The approach we take in this paper is to find a sequence of reductions refining each other with the guarantee that this sequence is of maximal possible length.

Definition 2.5 (Chain of lumpings). For an ODE system of the form $\mathbf{x}^{\prime}=\mathbf{f}(\mathbf{x})$, a sequence of linear transformations

$$
\mathbf{y}_{1}=\mathbf{x} L_{1}, \mathbf{y}_{2}=\mathbf{x} L_{2}, \ldots, \mathbf{y}_{\ell}=\mathbf{x} L_{\ell}
$$

where $L_{1} \in \mathbb{C}^{n \times m_{1}}, \ldots, L_{\ell} \in \mathbb{C}^{n \times m_{\ell}}$, is called a chain of lumpings if

1. $0<m_{1}<\ldots<m_{\ell}<n$;
2. $\mathbf{y}_{i}=\mathbf{x} L_{i}$ is a lumping of (2.1) for every $1 \leqslant i \leqslant \ell$;
3. for every $1<i \leqslant \ell$, there exists a matrix $A_{i}$ such that $L_{i-1}=L_{i} A_{i}$.

The latter means that the reductions given by $L_{1}, \ldots, L_{\ell}$ refine each other. Such a chain $\left(L_{1}, \ldots, L_{\ell}\right)$ will be called maximal if it not contained as a subsequence in any longer chain.

We will show (see Corollary 3.6) that all maximal chains are of the same length, so they are also the longest possible chains. Given a maximal chain of lumpings, a user can "zoom in/out" by going left/right along the chain depending on the desired tradeoff between the size of the reduced model and the amount of information retained. Thus, we can now formally state the main problem studied in this paper.

Main problem 2.6.
Given a system $\mathbf{x}^{\prime}=\mathbf{f}(\mathbf{x})$ with $\mathbf{f}$ being a vector of polynomials over $\mathbb{Q}$;
Compute a maximal chain of lumpings for the system.


Figure 1: Maximal chain of lumpings for (2.3) and the corresponding reductions

Example 2.7 (Maximal chain of lumpings for (2.3)). Figure 1 shows a chain of lumpings and the corresponding reductions for our example system (2.3). The blocks contain the reduced systems and the arrows are labeled with the transformations between the consecutive reductions (matrices $A_{i}$ in the terms of Definition 2.5). This chain of reductions includes our preceding Examples 2.2 and 2.3 as $\mathbf{y}_{1}$ and $\mathbf{y}_{3}$, respectively.

In this example we have the original dimension $n=5$ and the dimensions of the reductions $m_{1}=1, m_{2}=2, m_{3}=3, m_{4}=4$, so this chain is clearly maximal.
3. Algorithm. For finding a maximal chain of lumpings, we first use theory developed in [28] to reduce the problem to a problem about common invariant subspaces of a set of matrices (Subsection 3.1) and then solve the new problem using the structure theory of finite-dimensional algebras (Subsections 3.3 to 3.5 ). The overall algorithm is summarized in Subsection 3.6.
3.1. Reduction to the search for common invariant subspaces. Let $\mathbf{x}^{\prime}=\mathbf{f}(\mathbf{x})$ be an ODE system in variables $\mathbf{x}=\left(x_{1}, \ldots, x_{n}\right)$ and $\mathbf{f}$ being a row vector of polynomials $f_{1}, \ldots, f_{n} \in \mathbb{C}[\mathbf{x}]$. Let $J(\mathbf{x})$ be the Jacobian matrix of $\mathbf{f}$ with respect to $\mathbf{x}$. We denote the monomials in $\mathbf{x}$ appearing in $J(\mathbf{x})$ by $m_{1}(\mathbf{x}), \ldots, m_{N}(\mathbf{x})$. Then $J(\mathbf{x})$ can be written uniquely as

$$
J(\mathbf{x})=\left(\begin{array}{lll}
\nabla f_{1} & \ldots & \nabla f_{n}
\end{array}\right)=\sum_{i=1}^{N} J_{i} \cdot m_{i}(\mathbf{x}), \quad \text { where } \nabla g:=\left(\begin{array}{lll}
\frac{\partial g}{\partial x_{1}} & \ldots & \frac{\partial g}{\partial x_{n}} \tag{3.1}
\end{array}\right)^{T}
$$

and $J_{1}, \ldots, J_{N}$ are constant matrices.
Example 3.1. Consider the system

$$
x_{1}^{\prime}=x_{1}-2 x_{2}^{2}, \quad x_{2}^{\prime}=-x_{2}+x_{2}^{2}
$$

In this case the decomposition (3.1) will be

$$
J\left(x_{1}, x_{2}\right)=\left(\begin{array}{cc}
1 & 0 \\
-4 x_{2} & -1+2 x_{2}
\end{array}\right)=\left(\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right)+\left(\begin{array}{cc}
0 & 0 \\
-4 & 2
\end{array}\right) x_{2}
$$

Lemma 3.2. Using the notation above, linear transformation $\mathbf{y}=\mathbf{x} L$, where $L \in \mathbb{C}^{n \times m}$, is a lumping of $\mathbf{x}^{\prime}=\mathbf{f}(\mathbf{x})$ if and only if the column space of $L$ is invariant with respect to $J_{1}, \ldots, J_{N}$.

Proof. For the case $L \in \mathbb{R}^{n \times m}$, the statement follows from [28, Lemmas S.I. 1 and S.II.1]. The proof of [28, Lemmas S.I.1] remains correct after replacing $\mathbb{R}$ with $\mathbb{C}$, and the proof of [28, Lemmas S.II.1] will be correct for the case of $\mathbb{C}$ if the real inner products are replaced with the complex ones.

Corollary 3.3. A sequence of linear transformations $\mathbf{y}_{1}=\mathbf{x} L_{1}, \ldots, \mathbf{y}_{\ell}=\mathbf{x} L_{\ell}$ is a chain of lumpings if and only if the column spaces $V_{1}, V_{2}, \ldots, V_{\ell}$ of $L_{1}, \ldots, L_{\ell}$ satisfy

- $V_{i}$ is invariant with respect to $J_{1}, \ldots, J_{N}$ for every $1 \leqslant i \leqslant \ell$;
- $\{0\} \subsetneq V_{1} \subsetneq \ldots \subsetneq V_{\ell} \subsetneq \mathbb{C}^{n}$.

Furthermore, the chain of lumpings is maximal if and only if the chain $V_{1}, \ldots, V_{\ell}$ is not a subsequence of a chain of subspaces satisfying the two properties above.

In order to search for such chains of invariant subspaces, we will use theory of finite dimensional matrix algebras.

Definition 3.4 (Matrix algebra). Let $k$ be a field (e.g., $k=\mathbb{Q}, \mathbb{R}, \mathbb{C}$ ). A subspace $\mathcal{A} \subseteq k^{n \times n}$ of matrices is called an algebra if it is closed under multiplication and contains the identity matrix.

For a finite set $A_{1}, \ldots, A_{m} \in k^{n \times n}$, we denote the smallest algebra containing $A_{1}, \ldots, A_{m}$ by $\left\langle A_{1}, \ldots, A_{m}\right\rangle$. This algebra is equal to the span of all possible products these matrices.

For an ODE system $\mathbf{x}^{\prime}=\mathbf{f}(\mathbf{x})$ with $\mathbf{x}=\left(x_{1}, \ldots, x_{n}\right)$ and $f_{1}, f_{2}, \ldots, f_{n} \in \mathbb{C}[\mathbf{x}]$, we consider the coefficients $J_{1}, \ldots, J_{N}$ of the Jacobian matrix of $\mathbf{f}(\mathbf{x})$ written as a polynomial in $\mathbf{x}$ as in (3.1). We will call the algebra $\left\langle I_{n}, J_{1}, \ldots, J_{N}\right\rangle$ (where $I_{n}$ is the identity $n \times n$-matrix) the Jacobian algebra of the system $\mathbf{x}^{\prime}=\mathbf{f}(\mathbf{x})$.

## Example 3.5.

- Let $T_{n}$ be the set of all upper-triangular matrices in $k^{n \times n}$. Since the product of two upper-triangular matrices is upper-triangular again, $T_{n}$ is an algebra.
- Consider the system from Example 3.1. Its Jacobian algebra is

$$
\left\langle\left(\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right),\left(\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right),\left(\begin{array}{cc}
0 & 0 \\
-4 & 2
\end{array}\right)\right\rangle=\left\{M^{T} \mid M \in T_{2}\right\}
$$

Since a subspace $V \subset \mathbb{C}^{n}$ is invariant with respect to $J_{1}, \ldots, J_{N}$ if and only if it is invariant with respect to $\left\langle I_{n}, J_{1}, \ldots, J_{N}\right\rangle$, that is, invariant w.r.t. any element of the algebra, we will further focus on finding invariant subspaces of this Jacobian algebra. An immediate benefit is that we can use the Jordan-Hölder theorem [13, Theorem 1.5.1] to clarify our notion of the maximal chain of lumpings (Definition 2.5): the definition only requires that a maximal chain cannot be further refined, and this, in general, does not preclude the existence of longer chains. The following direct consequence of [13, Theorem 1.5.1] guarantees that a maximal chain indeed has the maximal possible length.

Corollary 3.6. For a given $O D E$ system $\mathbf{x}^{\prime}=\mathbf{f}(\mathbf{x})$, all maximal chains of lumpings have the same length.
3.2. Generating the algebra. For performing explicit computation with the Jacobain algebra $\left\langle I_{n}, J_{1}, \ldots, J_{N}\right\rangle$ (Definition 3.4), we will compute its basis. Algorithm 3.1 gives a simplified version of our approach, which is essentially [28, Algorithm 2] applied to matrices instead of vectors. Similarly to [28], we employ modular computation (cf. [28, Algorithm 3]) to avoid the intermediate expression swell and use sparse linear algebra.

Building upon this straightforward adaptation of the approach from [28], we significantly improve the performance by taking further advantage of the sparsity of the input and output. In applications, only a couple of $J_{1}, \ldots, J_{N}$ (typically, the constant term of (3.1)) are not so sparse, and the rest are extremely sparse; furthermore, the basis of the Jacobian algebra also can be often chosen to be very sparse. Hence, many of matrices $C$ from (Step 2)(b)i will be sparse as well. However, some of the matrices computed at the intermediate steps may be still quite dense slowing down the whole algorithm. We deal with the issue by temporarily deferring (Step 2)(b)ii for relatively dense matrices $C$ and then, once the outer loop exits signaling that $P$ is empty, we add each of the deferred matrices to $P$ and restart the iteration.

This way we ensure that we have generated enough sparse matrices in the algebra so that the reductions of the dense matrices will be more sparse now. Thanks to this optimization, Algorithm 3.1 is not a bottleneck in our computation which it was when we used the approach from [28] directly.

```
Algorithm 3.1 Finding a basis of matrix algebra (basic version)
Input a set of square matrices \(A_{1}, \ldots, A_{\ell} \in k^{n \times n}\);
Output a basis \(S\) of the smallest linear subspace \(\mathcal{A} \subseteq k^{n \times n}\) containing all possible products
        of \(A_{1}, \ldots, A_{\ell}\);
(Step 1) Set \(S\) to be any basis of the linear span of \(A_{1}, \ldots, A_{\ell}\) and let \(P:=S\).
(Step 2) While \(P \neq \varnothing\) do
(a) Take \(B\) to be an element of \(P\) and remove it from \(P\).
(b) For every \(A\) in \(\left\{A_{1}, \ldots, A_{\ell}\right\}\) do
            i. Compute \(C:=A B\) and reduce \(C\) w.r.t. \(S\) via Gaussian reduction.
            ii. If \(C \neq 0\), set \(S:=S \cup\{C\}\) and \(P:=P \cup\{C\}\).
```

(Step 3) Return $S$.
3.3. Search for invariant subspaces: algebraic preliminaries. For this section, we fix a ground field $k$ of characteristic zero. The cases we are mostly interested in are rational numbers $\mathbb{Q}$, algebraic numbers $\overline{\mathbb{Q}}$, and complex numbers $\mathbb{C}$.

Definition 3.7 (Radical of an algebra). Let $\mathcal{A} \subseteq k^{n \times n}$ be an algebra.

- $A$ subspace $\mathcal{I} \subseteq \mathcal{A}$ is called an ideal (resp., left ideal) if $A B, B A \in \mathcal{I}$ (resp., $A B \in \mathcal{I}$ ) for every $A \in \mathcal{A}$ and $B \in \mathcal{I}$.
- An ideal (resp., left ideal) $\mathcal{I} \subseteq \mathcal{A}$ is nilpotent if there exists $N$ such that the product of any $N$ elements of $\mathcal{I}$ is zero.
- The set of all elements $A \in \mathcal{A}$ such that the left ideal $\mathcal{A} \cdot A$ is nilpotent is called the radical of $\mathcal{A}$. It is a nilpotent ideal of $\mathcal{A}$ by [13, Theorems 3.1.6, 3.1.10].
Example 3.8. Let $T_{n}$ be the set of all upper-triangluar matrices in $k^{n \times n}$. Consider a subset $U_{n} \subset T_{n}$ of strictly upper-triangular matrices. One can easily verify that $U_{n}$ is an ideal and the product of any $n$ elements of $U_{n}$ is zero. Since, for every $A \in U_{n}$, we have $T_{n} \cdot A \subseteq U_{n}$, we deduce that $U_{n}$ is the radical of $T_{n}$.

Dixon's theorem [5, Theorem 11] implies that the radical of an algebra $\mathcal{A} \subseteq k^{n \times n}$ can be computed by finding the kernel of a square matrix of order $\operatorname{dim} \mathcal{A} \leqslant n^{2}$. The relevance of the notion of radical to our problem is demonstrated by the following lemma.

Lemma 3.9. Let $\mathcal{A} \subseteq k^{n \times n}$ be an algebra, and let $\mathcal{R} \subset \mathcal{A}$ be its radical. If $\mathcal{R} \neq\{0\}$, then the intersection $\bigcap_{R \in \mathcal{R}} \operatorname{Ker} R$ is nontrivial and is invariant w.r.t. $\mathcal{A}$.

Proof. Let $N$ be the smallest integer such that the product of any $N$ elements of $\mathcal{R}$ is zero. Then there exists $0 \neq M \in k^{n \times n}$ which is a product of $N-1$ elements of $\mathcal{R}$. Then, we have $R M=0$ for every $R \in \mathcal{R}$, so $V:=\bigcap_{R \in \mathcal{R}} \operatorname{Ker} R \supseteq \operatorname{Im} M$ is nontrivial.

Consider $v \in V, A \in \mathcal{A}$, and $R \in \mathcal{R}$. Since $R A \in \mathcal{R}$, we have $R A v=0$, so $A v \in \operatorname{Ker} R$. Thus, $V$ is invariant w.r.t. $\mathcal{A}$.

Example 3.10. Consider the system from Example 3.1. In Example 3.5, it was shown that the Jacobian algebra of this system is the set of lower triangular matrices. Similarly to Example 3.8 we find that the radical of this algebra is $\left(\begin{array}{ll}0 & 0 \\ \lambda & 0\end{array}\right)$. The common kernel of the radical is spanned by the second basis vector yielding the reduction $y^{\prime}=-y+y^{2}$ (with $y=x_{2}$ )

Definition 3.11 (Semisimple algebra). An algebra $\mathcal{A} \subseteq k^{n \times n}$ is called semisimple if its radical is zero.

We will use the following characterization of semisimple algebras.
Theorem 3.12 (Wedderburn-Artin, [13, Theorems 2.4.3 and 2.6.2]). Let $\mathcal{A} \subseteq k^{n \times n}$ be a semisimple algebra. Then there exist

1. algebras $\mathcal{A}_{1} \subseteq k^{n_{1} \times n_{1}}, \ldots, \mathcal{A}_{\ell} \subseteq k^{n_{\ell} \times n_{\ell}}$ such that $\mathcal{A}_{i}$ does not have a nontrivial proper invariant subspace in $k^{n_{i}}$ for every $1 \leqslant i \leqslant \ell$,
2. integers $m_{1}, \ldots, m_{\ell}$ such that $n_{1} m_{1}+\ldots+n_{\ell} m_{\ell}=n$,
3. a basis in $k^{n}$
such that, in this basis, we have

$$
\begin{equation*}
\mathcal{A}=\{\operatorname{Diag}(\underbrace{A_{1}, \ldots, A_{1}}_{m_{1} \text { times }}, \ldots, \underbrace{A_{\ell}, \ldots, A_{\ell}}_{m_{\ell} \text { times }}) \mid A_{1} \in \mathcal{A}_{1}, \ldots, A_{\ell} \in \mathcal{A}_{\ell}\} \tag{3.2}
\end{equation*}
$$

where $\operatorname{Diag}\left(B_{1}, \ldots, B_{N}\right)$ denotes the block-diagonal matrix with blocks $B_{1}, \ldots, B_{N}$.
Example 3.13. Consider the set of all matrices of the form

$$
\left(\begin{array}{cccc}
a & b & 0 & 0 \\
-b & a & 0 & 0 \\
0 & 0 & c & 0 \\
0 & 0 & 0 & c
\end{array}\right), \quad \text { where } a, b, c \in \mathbb{Q}
$$

This is a semisimple algebra in the form (3.2) with $\ell=2, m_{1}=1$, and $m_{2}=2$.
In the case $\ell=1$ and $m_{1}=1$ in the decomposition (3.2) from Theorem 3.12, there are no invariant subspaces in $k^{n}$ but there still may be invariant subspaces in $\bar{k}^{n}$ if $k \neq \bar{k}$, where $\bar{k}$ is the algebraic closure of field $k$. These subspaces can be found using the center of the algebra.

Definition 3.14 (Center/Centralizer). Let $\mathcal{A} \subseteq k^{n \times n}$ be an algebra.

- The center of $\mathcal{A}$ is the set of all $M \in \mathcal{A}$ such that $M A=A M$ for every $A \in \mathcal{A}$.
- The centralizer of $\mathcal{A}$ is the set of all $M \in k^{n \times n}$ such that $M A=A M$ for every $A \in \mathcal{A}$.

Since, for every fixed $A, A M=M A$ is a system of linear equations in the entries, the center and centralizer can be computed by solving a system of linear equations.

Lemma 3.15. Let $\mathcal{A} \subseteq k^{n \times n}$ be an algebra without nontrivial proper invariant subspaces in $k^{n}$. Let $\mathcal{C}$ be the center of $\mathcal{A}$. For every $C \in \mathcal{C}$, every eigenspace of $C$ is an invariant subspace of $\mathcal{A}$ in $\bar{k}^{n}$.

Proof. Let $V$ be an eigenspace of $C$ corresponding to the eigenvalue $\lambda$. Let $A \in \mathcal{A}$. Then $C(A v)=(C A) v=(A C) v=\lambda A v$.

Lemma 3.16. Let $\mathcal{A} \subseteq \mathbb{Q}^{n \times n}$ be a semisimple algebra. Let $M \in \mathcal{A}$ be a matrix such that the characteristic polynomial of $M$ is of the form $p(t)^{d}$, where $p(t)$ is $\mathbb{Q}$-irreducibe. Let $\mathcal{Z}$ and $\mathcal{C}$ be the center and centralizer of $\mathcal{A}$, respectively. Then the equality $\operatorname{dim} \mathcal{C}=d^{2} \operatorname{dim} \mathcal{Z}$ is equivalent to the fact that, in the Wedderburn-Artin decomposition (3.2) of $\mathcal{A}$, we have $\ell=1$ and $m_{1}=d$.

Proof. We consider the Wedderburn-Artin decomposition (3.2) of $\mathcal{A}$. For every $1 \leqslant i \leqslant \ell$, we denote the center of $\mathcal{A}_{i}$ by $\mathcal{Z}_{i}$. Then $\operatorname{dim} \mathcal{Z}=\operatorname{dim} \mathcal{Z}_{1}+\ldots+\mathcal{Z}_{\ell}$. The number of irreducible factors of a characteristic polynomial of any element of $\mathcal{A}$ will be at least $m_{1}+\ldots+m_{\ell}$, so $d \geqslant m_{1}+\ldots+m_{\ell}$. A direct computation using the Schur's lemma [13, Theorem 2.1.1] implies that the centralizer $\mathcal{C}$ of $\mathcal{A}$ is isomorphic to $\operatorname{Mat}_{m_{1}}\left(\mathcal{Z}_{1}\right) \times \ldots \times \operatorname{Mat}_{m_{\ell}}\left(\mathcal{Z}_{\ell}\right)$, where $\operatorname{Mat}_{m_{i}}\left(\mathcal{Z}_{i}\right)$ denotes the space of $m_{i} \times m_{i}$-block matrices with each block being an element of $\mathcal{Z}_{i}$ (cf. [13, Theorem 2.6.4]). Therefore

$$
\operatorname{dim} \mathcal{C}=m_{1}^{2} \operatorname{dim} \mathcal{Z}_{1}+m_{2}^{2} \operatorname{dim} \mathcal{Z}_{2}+\ldots+m_{\ell}^{2} \operatorname{dim} \mathcal{Z}_{\ell}
$$

Bounding the right-hand side, we can write

$$
\operatorname{dim} \mathcal{C} \leqslant\left(m_{1}+\ldots+m_{\ell}\right)^{2}\left(\operatorname{dim} \mathcal{Z}_{1}+\ldots+\operatorname{dim} \mathcal{Z}_{\ell}\right) \leqslant d^{2} \operatorname{dim} \mathcal{Z}
$$

Both inequalities will be equalities if and only if $\ell=1$ and $d=m_{1}$, and this proves the lemma.
3.4. Search for invariant subspaces: how to find one. In this subsection, we present Algorithm 3.2 for finding an invariant subspace if there is any. The rest of the subsection is devoted to justifying its correctness and termination, see Proposition 3.20.

Proposition 3.17. Let $\mathcal{A} \subseteq \mathbb{Q}^{n \times n}$ be a semisimple algebra such that there are no nontrivial proper $\mathcal{A}$-invariant subspaces in $\mathbb{Q}^{n}$. Let $M_{1}, \ldots, M_{N}$ be a linear basis of $\mathcal{A}$. Then the polynomial

$$
\operatorname{det}\left(x_{1} M_{1}+\ldots+x_{N} M_{N}\right) \in \mathbb{Q}\left[x_{1}, \ldots, x_{N}\right]
$$

is of the form $P^{d}$, where $P$ is irreducible over $\mathbb{Q}$.
Remark 3.18 (On the importance of being a basis). While the statement of Proposition 3.17 may sound quite natural, there situation is in fact quite subtle: if one replaces linear basis with a set of generators of $\mathcal{A}$ in the statement of the proposition, it will not longer be true [24, Theorem 1.2].

Proof of Proposition 3.17. By performing a change of coordinates over $\mathbb{Q}$, we will assume that $M_{1}$ is the identity matrix.

Let $\overline{\mathcal{A}}$ be the complexification of $\mathcal{A}$. By the Wedderburn-Artin theorem [13, Corollary 2.4.4], there exist $n_{1}, \ldots, n_{\ell}$ such that $N=n_{1}^{2}+\ldots+n_{\ell}^{2}$ and

$$
\begin{equation*}
\overline{\mathcal{A}} \cong \operatorname{Mat}_{n_{1}}(\mathbb{C}) \times \ldots \times \operatorname{Mat}_{n_{\ell}}(\mathbb{C}) \tag{3.3}
\end{equation*}
$$

```
Algorithm 3.2 Finding a nontrivial invariant subspace of an algebra
Input a basis \(B_{1}, \ldots, B_{N} \in \mathbb{Q}^{n \times n}\) of an algebra \(\mathcal{A} \subseteq \mathbb{Q}^{n \times n}\);
Output One of the following:
    - NO if there is no subspace in \(\overline{\mathbb{Q}}^{n}\) invariant w.r.t. \(\mathcal{A}\);
    - nontrivial proper subspace in \(\mathbb{Q}^{n}\) invariant w.r.t. \(\mathcal{A}\);
    - a maximal chain of subspaces in \(\overline{\mathbb{Q}}^{n}\) invariant w.r.t. \(\mathcal{A}\).
    Considering corner cases:
```

(Step 1) If $N=n^{2}$, return NO.
(Step 2) For an arbitrary nonzero vector $v$, consider a space $V$ spanned by $B_{1} v, \ldots, B_{N} v$. If
$\operatorname{dim} V<n$, return $V$.
Examining the radical:
(Step 3) Find a basis of the radical $\mathcal{R}$ of $\mathcal{A}$ (Definition 3.7) using Dixon's theorem [5, Theorem 11]
(Step 4) If $\operatorname{dim} \mathcal{R}>0$ compute the common kernel of the basis elements of $\mathcal{R}$ and return it (see Lemma 3.9).
Semisimple case:
(Step 5) Set $D:=1$.
(Step 6) Compute $M:=\sum_{i=1}^{N} a_{i} B_{i}$, where $a_{1}, \ldots, a_{N}$ are sampled independently and uniformly at random from $\{1,2, \ldots, D\}$.
(Step 7) If the characteristic polynomial of $M$ has at least two distinct $\mathbb{Q}$-irreducible factors (say, $p_{1}(t)$ and $p_{2}(t)$ ):
(a) Check the invariance of $\operatorname{Ker} p_{1}(M)$ w.r.t. $B_{1}, \ldots, B_{N}$.
(b) If it is invariant, return $\operatorname{Ker} p_{1}(M)$. Otherwise, set $D:=2 D$ and go to (Step 6).
(Step 8) Write the characteristic polynomial of $M$ as $p(t)^{d}$, where $p(t)$ is $\mathbb{Q}$-irreducible.
(Step 9) Compute the center $\mathcal{Z}$ and centralizer $\mathcal{C}$ of $\mathcal{A}$ (Definition 3.14).
(Step 10) If $\operatorname{dim} \mathcal{C}<d^{2} \operatorname{dim} \mathcal{Z}$, set $D:=2 D$ and go to (Step 6).
(Step 11) Let $C_{1}, \ldots, C_{s}$ be a basis of $\mathcal{C}$. Set $C:=\sum_{i=1}^{s} b_{i} C_{i}$, where $b_{1}, \ldots, b_{s}$ are sampled independently and uniformly at random from $\{1,2, \ldots, D\}$.
(Step 12) Compute $q(t)$, the minimal polynomial of $C$. If $q$ is $\mathbb{Q}$-reducible or $\operatorname{deg} q<d \operatorname{dim} \mathcal{Z}$, set $D:=2 D$ and go to (Step 6).
(Step 13) Let $V_{1}, \ldots, V_{\ell}$ (where $\ell=d \operatorname{dim} \mathcal{Z}$ ) be the eigenspaces of $C$.
(Step 14) Return $V_{1} \subset V_{1} \oplus V_{2} \subset \ldots \subset V_{1} \oplus V_{2} \oplus \ldots \oplus V_{\ell-1}$.

Then the complexification $\mathbb{C}^{n}$ of the original representation $\mathbb{Q}^{n}$ of $\mathcal{A}$ can be decomposed [13, Theorem 2.6.2] as

$$
\begin{equation*}
\mathbb{C}^{n}=k_{1} V_{1} \oplus k_{2} V_{2} \oplus \ldots \oplus k_{\ell} V_{\ell} \tag{3.4}
\end{equation*}
$$

where $V_{i} \cong \mathbb{C}^{n_{i}}$ is the unique irreducible representation of $\operatorname{Mat}_{n_{i}}(\mathbb{C})$. We denote the base change corresponding to (3.4) by $C \in \mathbb{C}^{n \times n}$. Then $C M C^{-1}$, where $M:=x_{1} M_{1}+\ldots+x_{N} M_{N}$, is block diagonal with the dimensions of blocks as in (3.4). Furthermore, there exists and invertible matrix $B \in \mathbb{C}^{N \times N}$ such that, for $\mathbf{y}:=B \mathbf{x}$, one has

$$
C M C^{-1}=\operatorname{diag}(\underbrace{Y_{1}, \ldots, Y_{1}}_{k_{1} \text { times }}, \ldots, \underbrace{Y_{\ell}, \ldots, Y_{\ell}}_{k_{\ell} \text { times }}),
$$

where $Y_{i}$ is a matrix with entries $y_{n_{1}+\ldots+n_{i-1}^{2}+1}, \ldots, y_{n_{1}^{2}+\ldots+n_{i}^{2}}$ for every $1 \leqslant i \leqslant \ell$.
Then we have

$$
\operatorname{det}(M)=\operatorname{det}\left(C M C^{-1}\right)=\operatorname{det}\left(Y_{1}\right)^{k_{1}} \ldots \operatorname{det}\left(Y_{\ell}\right)^{k_{\ell}}
$$

Furthermore, since $M_{1}$ is the identity, $\left.\operatorname{det}(M)\right|_{x_{1}=x_{1}+t}$ as a polynomial in $t$ is the characteristic polynomial of $-M$. Let $Q(\mathbf{x}):=\operatorname{det} Y_{1} \ldots \operatorname{det} Y_{\ell} \in \mathbb{Q}[\mathbf{x}]$. Then $\left.Q\right|_{x_{1}=x_{1}+t}$ as a polynomial in $t$ is the minimal polynomial of $-M$.

Since $\operatorname{det} Y_{i}$ is a determinant of a matrix with independent entries, it is irreducible over $\mathbb{C}$. Let $p(\mathbf{x})$ be a $\mathbb{Q}$-irreducible divisor of $\operatorname{det} M$. Then $p$ divides $Q$, so, by reordering $Y_{i}$ 's if necessary, we can assume that $p(\mathbf{x})=\operatorname{det} Y_{1} \ldots \operatorname{det} Y_{r}$ for $r \leqslant \ell$. Assume that $r<\ell$. Set $p_{0}(t):=p\left(x_{1}-t, x_{2}, \ldots, x_{N}\right)$ and consider $p_{0}(M)$. We will have

$$
C p_{0}(M) C^{-1}=\operatorname{diag}(\underbrace{0, \ldots, 0}_{k_{1}+\ldots+k_{r} \text { times }}, \underbrace{p_{0}\left(Y_{r+1}\right), \ldots, p_{0}\left(Y_{r+1}\right)}_{k_{r+1} \text { times }}, \ldots, \underbrace{p_{0}\left(Y_{r+1}\right), \ldots, p_{0}\left(Y_{r+1}\right)}_{k_{\ell} \text { times }}) .
$$

Since $p_{0}$ is coprime with the charpolys of $Y_{r+1}, \ldots, Y_{\ell}$, the matrices $p_{0}\left(Y_{r+1}\right), \ldots, p_{0}\left(Y_{\ell}\right)$ are nonsingular. Therefore, the kernel of $C p_{0}(M) C^{-1}$ is exactly the span of the first $k_{1}+\ldots+k_{r}$ basis vectors. Therefore, the kernel of $p_{0}(M)$ is the span of this many first columns of $C^{-1}$. Therefore, the kernel of $p_{0}(M)$ is $\overline{\mathcal{A}}$-invariant and is defined over $\mathbb{C}$. On the other hand, the entries of $p_{0}(M)$ belong to $\mathbb{Q}(\mathbf{x})$, so the kernel of $p_{0}(M)$ in fact is defined over $\mathbb{C} \cap \mathbb{Q}(\mathbf{x})=\mathbb{Q}$. Therefore, the kernel of $p_{0}(M)$ yields a nontrivial $\mathcal{A}$-invariant subspace of $\mathbb{Q}^{n}$ contradicting with the irreducibility of this representation. Therefore $p$ must be equal to $Q$ and, thus, det $M$ must be a power of $p$.

The proof of the proposition provides a way to find the degree of $\operatorname{deg} P$.
Corollary 3.19. In the notation of the proof (see (3.3)) of Proposition 3.17, $\operatorname{deg} P=n_{1}+$ $n_{2}+\ldots+n_{\ell}$.

Proposition 3.20. Algorithm 3.2 is correct and terminates with probability one.
Proof. We will first prove the correctness. If the algorithm returned on (Step 1), then $\mathcal{A}$ is the full matrix algebra, and does not have any nontrivial proper invariant subspace. If the algorithm returned on (Step 2), then the returned subspace is invariant by construction. If the algorithm returned on (Step 4), the returned subspace is nonzero and invariant due to Lemma 3.9.

It remains to consider the case when the algorithm returns after (Step 5). If the algorithm returned on (Step 7)ii., then the returned subspace is invariant by construction and is
nonzero because $p_{1}(t)$ divides the charpoly of $M$, so $p_{1}(M)$ is a singular matrix. Finally, consider the case when the algorithm returned on (Step 14). Consider the decomposition (3.2) from Theorem 3.12 for $\mathcal{A}$. If the algorithm reached (Step 11), it contains a matrix with the charpoly being $p(t)^{d}$ with $\mathbb{Q}$-irreducible $p(t)$ such that $\operatorname{dim} \mathcal{C}=d^{2} \operatorname{dim} \mathcal{Z}$. Lemma 3.16 implies that, in the decomposition (3.2), we have $\ell=1$ and $m_{1}=d$. Thus, the whole space $\mathbb{Q}^{n}$ can be written as $U_{1} \oplus U_{2} \oplus \ldots \oplus U_{d}$ such that each of $U_{i}$ 's is $\mathcal{A}$-invariant without proper nontrivial $\mathcal{A}$-invariant subspaces. [13, Corollary 2.2.4] implies that, over $\overline{\mathbb{Q}}$, each of $U_{i}$ 's can be decomposed as a direct sum of at most $\operatorname{dim} \mathcal{Z} \mathcal{A}$-invariant subspaces. Therefore, the whole $\overline{\mathbb{Q}}^{n}$ can be decomposed into at most $d \operatorname{dim} \mathcal{Z} \mathcal{A}$-invariant subspaces by [13, Theorem 2.6.2].

Lemma 3.15 implies that each of $V_{i}$ 's from (Step 13) is an invariant subspace w.r.t. $\mathcal{A}$. Since there are $d \operatorname{dim} \mathcal{Z}$ of them, each of $V_{i}$ 's does not contain nontrivial proper $\mathcal{A}$-invariant subspaces. Therefore, the chain $V_{1} \subset V_{1} \oplus V_{2} \subset \ldots \subset V_{1} \oplus V_{2} \oplus \ldots \oplus V_{\ell-1}$ returned at (Step 14) is maximal. This finished the proof of the correctness of the algorithm.

We will now prove that the algorithm terminates with probability one. Consider the decomposition (3.2) of $\mathcal{A}$ from Theorem 3.12. Consider variables $z_{1}, \ldots, z_{N}$ and a ma$\operatorname{trix} M_{0}:=\sum_{i=1}^{N} z_{i} B_{i}$. Then $M$ at (Step 6) is a specialization of $M_{0}$ at $z_{i}=a_{i}$. Let $P\left(z_{1}, \ldots, z_{N}, t\right)$ be the charpoly of $M_{0}$. Consider the decomposition (3.2) for $\mathcal{A}$. For every $1 \leqslant i \leqslant \ell$, we apply Proposition 3.17 to the block corresponding to $\mathcal{A}_{i}$ and obtain a $\mathbb{Q}$-irreducible $P_{i}$ and its power $d_{i}$. Thus, we obtain the following factorization for $M_{0}$

$$
P=P_{1}^{d_{1} m_{1}} P_{2}^{d_{2} m_{2}} \ldots P_{\ell}^{d_{\ell} m_{\ell}} .
$$

The characteristic polynomial of $M$ computed at (Step 6) is equal to $P\left(a_{1}, \ldots, a_{N}, t\right)$. Assume that $P_{i}\left(a_{1}, \ldots, a_{N}, t\right)$ is $\mathbb{Q}$-reducible for every $1 \leqslant i \leqslant s$ and these polynomials are distinct.

- Assume that $\ell>1$. Then $p_{1}(t)$ from (Step 7) will be equal to $P_{i}\left(a_{1}, \ldots, a_{N}, t\right)$ for some $i$. Then $\operatorname{Ker} p_{1}(M)$ will be the subspace corresponding to the $i$-th block in the decomposition (3.2). The subspace is invariant, so it will be returned on (Step 7)ii..
- Assume that $\ell=1$. We will study matrix $C$ similarly to the way we studied $M$ above. Let $y_{1}, \ldots, y_{s}$ be independent variables, and we define $C_{0}:=y_{1} C_{1}+\ldots+y_{s} C_{s}$. By the same argument as in the proof of Lemma 3.16, we have $\mathcal{C} \cong \operatorname{Mat}_{r}(\mathcal{Z})$ for some integer $r$. By [13, Proposition 2.3.4], algebra $\mathcal{C}$ is simple and every $\mathcal{C}$-module (in particular, our ambient space $\mathbb{Q}^{n}$ ) is a direct sum of isomorphic copies of the same $\mathcal{C}$-module. We apply Proposition 3.17 to this module and deduce that the characteristic polynomial of $C_{0}$ is of the form $Q\left(y_{1}, \ldots, y_{s}, t\right)^{h}$ for some integer $h$ and $\mathbb{Q}$-irreducible polynomial $Q$. Furthermore, $\operatorname{deg} Q=d \operatorname{dim} \mathcal{Z}$ by Corollary 3.19. Assume that $Q\left(b_{1}, \ldots, b_{s}, t\right)$ is $\mathbb{Q}$-irreducible. Then $Q\left(b_{1}, \ldots, b_{s}, t\right)$ will be the minimal polynomial of $C$, so this polynomial will not satisfy the condition of (Step 12) and, thus, the algorithm will terminate without going back to (Step 6).
Combining the two underlined assumptions in the text above, we see that the algorithm will return for a fixed value of $D$ if the following conditions hold:

1. $P_{i}\left(a_{1}, \ldots, a_{N}, t\right)$ is $\mathbb{Q}$-reducible for every $1 \leqslant i \leqslant s$ and these polynomials are all distinct;
2. $Q\left(b_{1}, \ldots, b_{s}, t\right)$ is $\mathbb{Q}$-irreducible.
[11, Theorem 2.1] implies that there exists constants $C_{0}, C_{1}$ such that the probability of any of $P_{i}\left(a_{1}, \ldots, a_{N}, t\right)$ 's and $Q\left(b_{1}, \ldots, b_{s}, t\right)$ being $\mathbb{Q}$-reducible is less that $\frac{C_{1}}{\sqrt[3]{D}}$ if $D>C_{0}$. Fur-
thermore, the Schwartz-Zippel lemma [38, Proposition 98] implies that there exists a constant $C_{2}$ such that the probability of any of $P_{i}\left(a_{1}, \ldots, a_{N}, t\right)$ 's being equal does not exceed $\frac{C_{2}}{D}$. Therefore, for $D>C_{0}$, the probability that $D$ will be updated is at most $\frac{C_{1}}{\sqrt[3]{D}}+\frac{C_{2}}{D}$. This number will eventually become less than 0.99 , so the probability of non-termination will be bounded by $0.99 \cdot 0.99 \cdot \ldots=0$.
3.5. Search for invariant subspaces: how to find a chain. In this section, we describe how to use Algorithm 3.2 in a recursive manner to find a maximal chain of invariant subspaces in $\overline{\mathbb{Q}}$ w.r.t. the Jacobian algebra $\mathcal{A} \subset \mathbb{Q}^{n \times n}$ of an ODE system. We will denote a basis of $\mathcal{A}$ by $B_{1}, \ldots, B_{N}$

In the cases when Algorithm 3.2 applied to $B_{1}, \ldots, B_{N}$ returned NO or a maximal chain of invariant subspaces, we are done. Therefore, we consider the case when Algorithm 3.2 returns a single invariant subspace $V \subset \mathbb{Q}^{n}$. In this case, we consider two subproblems:

1. Restriction. Since $V$ is invariant w.r.t. $B_{1}, \ldots, B_{N}$, there are well-defined restrictions $\left.B_{1}\right|_{V}, \ldots,\left.B_{N}\right|_{V}$. We fix a basis in $V$ and will denote the matrix representations for these restricted operators also by $B_{1}^{*}, \ldots, B_{N}^{*}$.
2. Quotients. Consider the quotient space $\mathbb{Q}^{n} / V$ and the quotient map $\pi: \mathbb{Q}^{n} \rightarrow \mathbb{Q}^{n} / V$ (see $[3,3.83,3.88]$ ). Since $V$ is invariant w.r.t. $B_{1}, \ldots, B_{N}$, we can consider the quotient operators [3, 5.14] $B_{1} / V, \ldots, B_{N} / V$, we denote their matrix representations by $B_{1}^{\circ}, \ldots, B_{N}^{\circ}$. Note that, for every their common invariant subspace $U \subset \mathbb{Q}^{n} / V$, the subspace $\pi^{-1}(U) \subset \mathbb{Q}^{n}$ is invariant w.r.t. $B_{1}, \ldots, B_{N}$.
Note that the aforementioned matrix representations can be computed solving linear systems in $n$ variables. Thus, we can work recursively with algebras $\left\langle B_{1}^{*}, \ldots, B_{N}^{*}\right\rangle$ on $V$ and $\left\langle B_{1}^{\circ}, \ldots, B_{N}^{\circ}\right\rangle$ on $\mathbb{Q}^{n} / V$. If the resulting maximal chains of invariant subspaces are

$$
0 \subsetneq V_{1} \subsetneq \ldots \subsetneq V_{s} \subsetneq V \quad \text { and } \quad 0 \subsetneq U_{1} \subsetneq \ldots \subsetneq U_{r} \subsetneq \mathbb{Q}^{n} / V,
$$

then we can return the following maximal chain of invariant subspaces for $B_{1}, \ldots, B_{N}$

$$
0 \subsetneq V_{1} \subsetneq \ldots \subsetneq V_{s} \subsetneq V \subsetneq \pi^{-1}\left(U_{1}\right) \subsetneq \ldots \subsetneq \pi^{-1}\left(U_{r}\right) \subsetneq \mathbb{Q}^{n} .
$$

3.6. Putting everything together. In this section we collect the subroutines from the preceding sections into the complete algorithm for finding a maximal chain of lumpings.

```
Algorithm 3.3 Finding a maximal chain of lumpings
Input an ODE system \(\mathbf{x}^{\prime}=\mathbf{f}(\mathbf{x})\) with \(\mathbf{x}=\left(x_{1}, \ldots, x_{n}\right)\) and \(\mathbf{f}=\left(f_{1}, \ldots, f_{n}\right) \in \mathbb{Q}[\mathbf{x}]^{n}\);
Output a maximal chain of lumpings (see Definition 2.5 and Example 2.7);
(Step 1) Compute the Jacobian \(J(\mathbf{x})\) of \(\mathbf{f}\) and the matrices \(J_{1}, \ldots, J_{\ell} \in \mathbb{Q}^{n \times n}\) from its de- composition as in (3.1).
(Step 2) Use Algorithm 3.1 to compute the basis \(B_{1}, \ldots, B_{N}\) of the Jacobian algebra \(\mathcal{A}=\) \(\left\langle I_{n}, J_{1}, \ldots, J_{\ell}\right\rangle\) of the system.
```

(Step 3) Apply Algorithm 3.2 in a recursive way as decribed in Subsection 3.5 to compute a maximal chain $V_{1} \subsetneq \ldots \subsetneq V_{s}$ of subspaces in $\overline{\mathbb{Q}}^{n}$ invariant w.r.t. $\mathcal{A}$.
(Step 4) For each $1 \leqslant i \leqslant s$, find a matrix $L_{i}$ with the columns being a basis of $V_{i}$.
(Step 5) Return $L_{1}, \ldots, L_{s}$.
4. Implementation and performance. We have implemented Algorithm 3.3 (and all the algorithms it relies on) in Julia language [4] as a part of ExactODEReduction.jl package. The package together with relevant resources to replicate our results is freely available at https://github.com/x3042/ExactODEReduction.j1
We use libraries AbstractAlgebra.jl and Nemo.jl [18]. Internally, this results in using FLINT [21] and Calcium [23] (for complex number arithmetic). We use a version of the code from [30] to improve interpretability of the computed lumpings. Additionaly, during the development stage, various components of the package were profiled on collections of sparse matrices from the SuiteSparse dataset [12]. Our implementation accepts models typed manually or from the files in the ERODE *.ode format [6, Section 3.2]. We provide documentation, installation instructions, and usage examples.

We will demonstrate the performance of our implementation on a set of benchmarks ${ }^{2}$. We use benchmarks from the BioModels database [27] collected in [29] of dimensions ranging from 4 to 133. We run Algorithm 3.3 over rationals on each of the models. Table 1 contains benchmark results aggregated by models' dimension. For each range, we report:

- the number of models considered;
- the (average) length of a chain of reductions found;
- the (average) number of nonequivalent reductions, where equivalence is taken up to adding states with constant dynamics. We have chosen to report this because we think is it a reasonable first approximation to the number of "interesting" reductions;
- the (minimum, average, maximum) elapsed runtime of our implementation;

| Models info <br> Dimension |  | Reductions |  |  | Runtime (sec.) |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $2-9$ | 44 | 4.02 | 1.39 | 0.0 s | 0.6 s | 0.66 s |  |
| $10-19$ | 41 | 8.15 | 2.61 | 0.01 s | 0.21 s | 1.46 s |  |
| $20-29$ | 46 | 9.65 | 2.13 | 0.08 s | 0.44 s | 1.48 s |  |
| $30-39$ | 17 | 19.41 | 2.71 | 0.33 s | 1.74 s | 5.91 s |  |
| $40-59$ | 25 | 29.08 | 6.08 | 0.78 s | 4.58 s | 26.71 s |  |
| $60-79$ | 20 | 37.25 | 6.95 | 7.7 s | 34.57 s | 102.92 s |  |
| $80-99$ | 11 | 42.91 | 7.09 | 24.46 s | 96.38 s | 497.26 s |  |
| $100-133$ | 4 | 89.0 | 21.5 | 75.15 s | 202.52 s | 312.02 s |  |

Table 1: Benchmark results aggregated by model dimension

The timings were produced on a laptop with 2 cores 1.60 GHz each and $8 \mathrm{~Gb} \mathrm{RAM}^{3}$. We would like to note that out of the 208 models considered, at least one reduction was found in 202 models, and 154 of them admit a non-constant reduction.

The timings in the table do not include the cost of the positivization step [30], which is

[^2]optional. Here, our algorithm uses the Polymake [2] library. With the positivization step, the running time increases no more than by a factor of two in most instances, and usually the increase is indistinguishable at all ${ }^{4}$. In the earlier versions of the implementation of Algorithm 3.3, computing the algebra basis on (Step 2) had often been a clear bottleneck on our dataset. With the modifications to the Algorithm 3.1 as described in Subsection 3.2, currently, the most time-consuming steps are the restriction and quotienting procedures applied on (Step 3) of Algorithm 3.3. Solving a number of linear systems to find the matrix representations of restricted and quotient operators is a clear bottleneck here.

## 5. Case studies.

5.1. Inactivation of factor Va. We will consider a model from [22] which appears in the BioModels database [27] as BIOMD0000000365. Factor V is a protein involved in the process of coagulation (transforming blood from liquid to gel), and thus is closely related to blood vessel repair and thrombosis. In particular, it can assist in activating protein anticoagulant protein C. The activated factor V, factor Va, can no longer do this. A model describing deactivation of Va by means of activated protein C (APC) was proposed and studied in [22].

Factor Va consists of the heavy chain (HC) and light chain (LC), and the binding of APC happens through the light chain. The model consists of the following species

- Factor Va and its versions $\mathrm{Va}_{3}, \mathrm{Va}_{5}, \mathrm{Va}_{6}, \mathrm{Va}_{53}, \mathrm{Va}_{56}, \mathrm{Va}_{36}$, and $\mathrm{Va}_{536}$;
- LC, HC , and the versions of the latter $\left(\mathrm{HC}_{3}, \mathrm{HC}_{5}\right.$, etc) corresponding to the versions of Va;
- the $\mathrm{A}_{1}$ domain of factor $\mathrm{Va}, \mathrm{Va}_{L C \cdot A 1}$ and versions of the $\mathrm{A}_{2}$ domain such as $\mathrm{Va}_{A 3}$, $\mathrm{Va}_{A 53}$, etc.
- APC, complexes formed by it and LC/Va (such as APC•Va3).

In total, the model contains 30 variables and 9 parameters. Our code finds a maximal chain of lumpings of length 14 in under 5 second on a laptop. The smallest reduction with nonzero dynamics has dimension three and involves two parameters (similar to the one in Example 2.7):

$$
\left\{\begin{array}{l}
y_{1}^{\prime}=-k_{1} y_{1} y_{2}+k_{2} y_{3},  \tag{5.1}\\
y_{2}^{\prime}=-k_{1} y_{1} y_{2}+k_{2} y_{3}, \\
y_{3}^{\prime}=k_{1} y_{1} y_{2}-k_{2} y_{3} .
\end{array}\right.
$$

The macro-variables are

$$
\begin{aligned}
y_{1} & =[\mathrm{APC}], \\
y_{2} & =[\mathrm{LC}]+[\mathrm{Va}]+\left[\mathrm{Va}_{3}\right]+\left[\mathrm{Va}_{36}\right]+\left[\mathrm{Va}_{5}\right]+\left[\mathrm{Va}_{53}\right]+\left[\mathrm{Va}_{536}\right]+\left[\mathrm{Va}_{56}\right]+\left[\mathrm{Va}_{\mathrm{LC}} \cdot A_{1}\right], \\
y_{3} & =[\mathrm{LC} \cdot \mathrm{APC}]+[\mathrm{Va} \cdot \mathrm{APC}]+\left[\mathrm{Va}_{3} \cdot \mathrm{APC}\right]+[\mathrm{Va} 36 \cdot \mathrm{APC}]+\left[\mathrm{Va}_{5} \cdot \mathrm{APC}\right] \\
& +[\mathrm{Va} 53 \cdot \mathrm{APC}]+\left[\mathrm{Va}_{536} \cdot \mathrm{APC}\right]+[\mathrm{Va} 56 \cdot \mathrm{APC}]+\left[\mathrm{Va} \mathrm{VC}_{5} \cdot A_{1} \cdot \mathrm{APC}\right] .
\end{aligned}
$$

Variable $y_{2}$ (resp., $y_{3}$ ) can be described as the total concentration of the light chains without (resp., with) bound APC. Therefore, the reduction (5.1) focuses on the process of binding/unbinding of APC to the light chains, and it turns out that the other processes such as

[^3]reactions between the heavy and light chains become irrelevant and, in particular, the $\mathrm{HC}_{n}$ species do not appear in the macro-variables at all.

(a) States of the original model appearing in $y_{2}$

(b) Macro-variable $y_{2}$ of the reduced model

(c) Macro-variables $y_{1}, y_{3}$ of the reduced model

Figure 2: Numerical simulation for the model from [22] and its reduction using the initial conditions and parameter values from [22]

From numerical perspective ${ }^{5}$, the reduction (5.1) can be interpreted as "exact timescale separation" since the dynamics of the macro-variables turns out to be transient compared to the dynamics of the original system. More precisely, the original system was studied in [22] and has nontrivial dynamics on the timespan of 1200 second. In particular, this is the case for the variables contributing to the macro-variable $y_{2}$, see Figure 2a. On the other hand, as Figures 2 b and 2c show, the macro-variables $y_{1}, y_{2}, y_{3}$ have much faster dynamics and reach the steady state after less than one second.

[^4]

Figure 3: The relevant chemical species and dependencies between them
5.2. Model of cell death. In this subsection, we consider a model designed in [34] in order to study the sensitivity of the apoptosis (programmed cell death) to the TNF (tumor necrosis factor) stimulation. The overall model involves 47 chemical species and numerous interactions between them schematically described in [34, Figure 1]. Our code produces a maximal chain of lumpings of length 23 ( 16 out of them with nonconstant dynamics).

We will consider the nonconstant reduction of the smallest dimension. It involves two proteins, A20 and FLIP, whose concentrations depend on the concentrations of the corresponding mRNAs, A20_mRNA and FLIP_mRNA. The concentrations of these mRNAs are governed by the concentrations of nuclear $\mathrm{NF}-\kappa \mathrm{B}\left(\mathrm{NFkB} \_\mathrm{N}\right)$. The latter depends (directly or indirectly) on many other species including the aforementioned protein A20.

These species and relations between them are summarized on Figure 3, and the corresponding differential equations are:

$$
\begin{array}{ll}
{[A 20]^{\prime}=k_{1}\left[A 20 \_m R N A\right]+k_{2},} & {\left[A 20 \_m R N A\right]^{\prime}=k_{5}\left[N F \kappa B \_N\right],} \\
{[F L I P]^{\prime}=k_{3}\left[F L I P_{m} R N A\right]+k_{4},} & {\left[F L I P_{\_} m R N A\right]^{\prime}=k_{6}\left[N F \kappa B \_N\right],}
\end{array}
$$

where $k_{1}, \ldots, k_{6}$ are numeric parameters. Our code finds a three-dimensional reduction which can be straightforwardly simplified further a two-dimensional with the following macrovariables $y_{1}, y_{2}$ and the reduced system:

$$
\left\{\begin{array} { l } 
{ y _ { 1 } = \frac { k _ { 6 } } { k _ { 1 } } [ A 2 0 ] - \frac { k _ { 5 } } { k _ { 3 } } [ F L I P ] , } \\
{ y _ { 2 } = k _ { 6 } [ A 2 0 \_ m R N A ] - k _ { 5 } [ F L I P \_ m R N A ] }
\end{array} \Longrightarrow \left\{\begin{array}{l}
y_{1}^{\prime}=y_{2}+\frac{k_{2} k_{6}}{k_{1}}-\frac{k_{4} k_{5}}{k_{3}}, \\
y_{2}^{\prime}=0
\end{array}\right.\right.
$$

So the idea is that, although both A20 and FLIP are involved in a complex reaction network, one can, by eliminating the dependence on $\mathrm{NF} \kappa \mathrm{B}$, find a linear combination of them satisfying a simple system of differential equations which can be explicitly solved. Such explicit relations on the states can be, for example, combined with the differential inequalities method in order to obtain tighter reachability bounds [35].

By going further along the chain of the reductions one can include gradually more species into the reduced model, for example, a combination of the RIP protein and the transitional receptor can be included in a similar fashion.
6. Conclusions. We have presented a new algorithm which takes as input a system of ODEs and produces a longest possible chain of exact linear reductions of the system such that each reduction in the chain is a refinement of the previous one. This specification is more
general compared to the existing tools as it does not put any restriction on the new variables other than being the linear combinations of the original ones and it does not require any initial observable/guess.

We provided a publicly available implementation in Julia. Our code is able to analyze models of dimension over a hundred in a couple of minutes using commodity hardware. We have also demonstrated its applicability to models arising in life sciences.

Since the produced reductions are exact, our tool can be used for formal verification and as a preprocessing for approximate reduction techniques. While exactness is thus an important feature, it can also be viewed as a limitation since some models have only a few exact reductions (if any). Therefore, one intriguing direction for future research is to produce a "relaxed" version of our algorithm to find approximate lumpings together with rigorous error bounds. For existing results on approximate lumping, see $[37,19]$ and references therein. Interestingly, the core linear algebraic problem of our algorithm, finding common invariant subspaces, has been recently studied from the perspective of approximate but rigorous computation in [20, 10] motivated by factoring linear differential operators. We expect the ideas from these papers to be useful in our context as well.

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[^1]:    ${ }^{1}$ For the reduction in this example, it is sufficient that the rates of $E S \rightarrow E^{+}+S$ and $E \rightarrow E^{+}$coincide

[^2]:    ${ }^{2}$ Models are available at https://github.com/x3042/ExactODEReduction.jl/tree/main/data/ODEs, commit hash 678d32c5bbc8beedc9e22b673238cde0ec673a46.
    ${ }^{3}$ For the overall table, we refer to https://github.com/x3042/ExactODEReduction.jl/blob/main/ benchmark/biomodels_benchmark_results.md, commit hash 23c9f532aa316cbef59a8e3e6be04156a3d9c3eb.

[^3]:    ${ }^{4}$ One notable exception are models that admit large reductions with large coefficients. For example, model BIOMD0000000153 of dimension 76 has 22 nontrivial reductions of dimensions 55 and more, and applying the positivization routine increases the total runtime from 40 s to 1240 s .

[^4]:    ${ }^{5}$ All numerical simulations in this paper have been done using ModelingToolkit [26] and DifferentialEquations.jl [31]

