# Invariant Grids: Method of Complexity Reduction in Reaction Networks 

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## Key Words

Kinetics • Model reduction • Grids • Invariant manifold • Entropy • Non-linear dynamics - Mathematical modelling - Numerical methods


#### Abstract

Complexity in the description of big chemical reaction networks has both structural (number of species and reactions) and temporal (very different reaction rates) aspects. A consistent way to make model reduction is to construct the invariant manifold which describes the asymptotic system behaviour. In this paper we present a discrete analog of this object: an invariant grid. The invariant grid is introduced independently from the invariant manifold notion and can serve to represent the dynamic system behaviour as well as to approximate the invariant manifold after refinement. The method is designed for pure dissipative systems and widely uses their thermodynamic properties but allows also generalizations for some classes of open systems. The method is illustrated by two examples: the simplest catalytic reaction (Michaelis-Menten mechanism) and the hydrogen oxidation.


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## Simplexus

A large chemical reaction network can serve as a good model for imitating and predicting behaviour in various kinds of complex systems with interacting components, according to Gorban, Karlin, and Zinovyev. Complexity in the description of such a reaction network has both structural (number of chemicals and the reactions in which they take part) and temporal (different reaction rates) properties. For the network of biochemical reactions that take place in a living cell, the picture becomes very complicated, so a way of simplifying the network is needed. One consistent way of modelling a complex network is to construct its invariant manifold which describes its characteristics in multidimensional space of species concentrations. To glean information from the model, one then observes its behaviour projected onto the manifold.

In the present paper, the researchers have created a discrete analog of this the invariant manifold: an invariant grid. The invariant grid can represent dynamic system behaviour as well as provide an approximation of the reaction networks invariant manifold. The team explains that their approach has been developed for pure dissipative systems, in which the entropy grows monotonically and so utilizes thermodynamic properties of such systems. On the other hand, it is constructed so that it also allows them to generalize for numerous classes of open systems. For illustration, the team has focused on the oxidation of hydrogen and the simplest possible catalytic reaction, one that follows a conventional Michaelis-Menten mechanism.
The ultimate aim of this research is to provide a methodology for reducing the complexity of complex reactions and other systems and so allow new insights to be gained. Since all rate constants and complete reaction laws are rarely available for many complex reaction networks, their approach, which reduces the number of sys-

## 1 Introduction

Reaction networks serve as a good model to imitate and predict behaviour of complex systems of interacting components. Modern research is faced with a constantly increasing complexity of the systems under study: a good example is the fact that nowadays one can observe a boom in connection with studies of biochemical processes in a living cell [for recent overviews, see 1,2$]$. There is no need to underline an emerging need for methods of reducing the complexity of system description and behaviour.

The complexity in modelling big chemical reaction networks has both structural (number of species and reactions) and temporal (very different reaction rates) aspects (see fig. 1). In general, it is not possible to disregard the temporal organization of the network when one wants to create a realistic system model. of course, the rate constants and reaction laws are rarely available completely. This makes extremely desirable the development of methods allowing to reduce the number of system parameters as well as methods for qualitative analysis of chemical reaction networks [2].

The idea of model reduction with respect to slow motion extraction can be in-
troduced as follows: we have a system of ordinary differential equations describing time evolution of $n$ species concentrations (or masses) in time:

$$
\begin{equation*}
\frac{d \boldsymbol{x}}{d t}=J(\boldsymbol{x}) . \tag{1}
\end{equation*}
$$

Every particular state of the system corresponds to a point in the phase space $U$ and the system dynamics is determined by the vector field $J(x), x \in U$. We construct new (reduced) dynamics

$$
\begin{equation*}
\frac{d \boldsymbol{y}}{d t}=J^{\prime}(\boldsymbol{y}), \tag{2}
\end{equation*}
$$

where $y_{i}, i=1 \ldots m, m \ll n$ is a new set of variables corresponding to the slow dynamics of the initial system (1). By analogy with statistical physics it corresponds to the 'macroscopic' description of the chemical system (we observe only effects of slow system changes, comparable in time scale with characteristic times of experimental measurements) as opposed to 'microscopic' variables $x_{i}$. The reduced system dynamics exists on an $m$-dimensional manifold (surface) $\Omega$ embedded in the $n$-dimensional phase space and defined by functions $x_{i}=x_{i}\left(y_{1}, \ldots, y_{m}\right)$.

A consistent way of model reduction is to construct a positively invariant slow


Fig. 1. Graphical representation of two model systems considered as examples in this paper: Michaelis-Menten mechanism (a) and hydrogen burning model with 6 variables (b). Here circles represent chemical species, squares represent chemical reactions. Line thickness reflects direct reaction rate constants, a thicker line corresponds to a slower reaction (in a logarithmic scale). All reactions here are governed by mass action law and supposed to be reversible.
tem parameters required to be specified or measured with precision, will facilitate the qualitative analysis of chemical reaction networks.

To begin with, the team considered an approach to model reduction with respect to slow motion extraction in which a system of ordinary differential equations describing time evolution of a number of species of variable concentration (or mass) is sampled on a slow timescale, so, by analogy with statistical physics, a macroscopic description is obtained for a chemical system without considering fully detailed dynamics of microscopic variables. The dynamics of this reduced system can then be represented on a multidimensional manifold within a phase space of dimensions equal to the number of different chemical species involved. The pattern of positively invariant manifold is formed by slow motion segments of the system individual trajectories. The goal is to filter out the manifold and so provide a clear picture of how a dissipative dynamical system approaches its equilibrium.

In order to do this for a complex chemical reaction networks a computationally effective method is needed to build the invariant manifold in the first place. Thus constructing a surface of relatively low dimensionality can be reduced to a gridbased method. In this paper, Gorban, Karlin, and Zinovyev detail this method of invariant grids (MIG) approach. They point out that a grid could be refined repeatedly to bring it closer to the idea of the invariant manifold but by defining an invariant grid as an object independent of the manifold. An invariant grid, they explain further, is an undirected graph consisting of a set of nodes and the connections between them. The graph can be represented in a low-dimensional space with reduced coordinates. In this form, it is simply a finite lattice (regular and rectangular or hexagonal) but is simultaneously embedded in the phase space of the reaction system so that each node corresponds to a combination of
manifold $\Omega_{i n v}$, such that if an individual trajectory of the system (1) has started on $\Omega_{i n v}$, it does not leave $\Omega_{i n v}$ anymore, i.e. the vector field $J(x)$ in the points of the manifold is tangent to it (fig. 2a). The 'ideal' picture of the reduced description we have in mind is as follows: a typical phase trajectory, $x(t)$, where $t$ is the time, and $x$ is an element of the phase space, consists of two pronounced segments. The first segment connects the beginning of the trajectory, $x(0)$, with a certain point, $x\left(t_{1}\right)$, on the manifold $\Omega_{i n v}$ (rigorously speaking, we should think of $x\left(t_{1}\right)$ not on $\Omega_{i n v}$ but in a small neighbourhood of $\Omega_{i n v}$ but this is not essential for the ideal picture). The second segment belongs to $\Omega_{i n v}$. Thus, the manifolds appearing in our ideal picture are 'patterns' formed by the segments of individual trajectories, and the goal of the reduced description is to 'filter out' this manifold (fig. 2a).

Usually the construction of invariant manifold in the explicit form is difficult. Most of the time one deals with its approx-
imation constructed using a method [for overview, see 3-8]. It is formally possible to induce new dynamics on any given manifold $\Omega$, not necessarily invariant, if one introduces a projector operator $P$ of the vector field on the tangent bundle of the manifold $\Omega: P J(x \in \Omega) \in T_{x} \Omega$. By definition, the manifold $\Omega$ is invariant with respect to the vector field $J$ if and only if the following equality is true for each $x \in \Omega$ :

$$
\begin{equation*}
[1-P] J(x)=0 \tag{3}
\end{equation*}
$$

where projector $P$ depends on the point $x$ and on the manifold $\Omega$ in the vicinity of $x$. This equation is a differential equation for functions that define the manifold $\Omega$. The Newton method and the relaxation method, both iterative, were proposed to find a sequence of corrections to some initial approximation $\Omega$, in such a way that every next approximation has less invariance defect $[1-P] J(x)[$ see 5]. These corrections can be performed analytically in some cases.


Fig. 2. Main geometrical structures of model reduction: $U$ is the phase space, $J(x)$ is the vector field of the system under consideration: $d x / d t=J(x), \Omega$ is an ansatz manifold, $W$ is the space of macroscopic variables (coordinates on the manifold), the map $F: W \rightarrow U$ maps any point $y \in W$ into the corresponding point $x=F(y)$ on the manifold $\Omega, T x$ is the tangent space to the manifold $\Omega$ at the point $x, P J(x)$ is the projection of the vector $J(x)$ onto tangent space $T x$, the vector field $d y / d t$ describes the induced dynamics on the space of parameters, $\Delta=(1-P) J(x)$ is the defect of invariance, the affine subspace $x+\operatorname{ker} P$ is the plain of fast motions, and $\Delta \in$ $\operatorname{ker} P$. a Here $\Omega_{\text {inv }}$ is an invariant manifold (all $J\left(x \in \Omega_{\text {inv }}\right)$ are tangent to $\Omega_{\text {inv }}$ ) and a possible dynamics is shown in its vicinity. b Here $\Omega$ is some manifold approximating the invariant manifold ( $J(x \in \Omega$ ) is not necessarily tangent to $\Omega$ ), one can use operator $P$ to derive new dynamics (2).
chemical species concentrations. The connectivity of such a graph is needed to calculate a tangent space in every node using differentiation operators.
Gorban, Karlin, and Zinovyev here propose two algorithms, both iterative - a Newtonian type algorithm and a relaxation method that can be used to produce such a graph. They have demonstrated proof of principle using dissipative reaction systems so that their thermodynamic properties can be used to clearly define the metrics of the phase space and so allow them to perform geometrical calculations on their graph. Low-dimensional invariant manifolds do exist for other types of systems presenting more complicated dynamic behaviour but most of the physically significant models include non-dissipative components and thus can be analyzed with use of thermodynamics.

The researchers applied the MIG approach first to a simple two-step catalytic reaction in which the reaction of two reactants is promoted by a catalyst to produce a product and a by-product. The researchers assume a standard Michaelis-Menten mechanism for the chemical changes that take place. Their second example is a model of the hydrogen burning (oxidation) reaction. This is more complicated as there are six concentrations to consider (molecular hydrogen, oxygen, and water, and three radicals, $\mathrm{H}, \mathrm{O}$, and OH ). The resulting invariant grids provide much information about the character of system slow dynamics from different perspectives, how it proceeds depending on initial concentrations and what effect has viewing it close or far away from the equilibrium. The ultimate achievement is that a complex system becomes far easier to visualize using the gridbased approach than attempting to appreciate a detailed multidimensional vector field.

The team emphasize that their approach allows one to visualize simultaneously many different scenarios of system behaviour, together with different characteristics

For the case of a complex chemical reaction network, one has to develop a computationally effective method of invariant manifold construction. If one constructs a surface of a relatively low dimension, gridbased manifold representations become a relevant option [8]. In this paper we present such an approach named method of invariant grids (MIG).On one hand, grid representation can be refined and converge more and more closely to the invariant manifold. On the other, we define invariant grid as an object independent of the manifold itself. Thus, it can be used independently: for example, for visualization of the global system dynamics as will be shown at the end of this paper.

Invariant grid is an undirected graph which consists of a set of nodes and connections between them. The graph can be represented in two spaces: in the low-dimensional space of the internal (reduced) coordinates where it forms a finite lattice (usually, regular and rectangular or hexagonal), and, simultaneously, it is embedded in the phase space $U$; thus every node corresponds to a species concentration vector $x$. Using connectivity of the graph, one can introduce differentiation operators and calculate the tangent vectors and define the projector operator in every node. This is the only place where the connectivity of the graph is used. The node positions in $U$ are optimized such that the invariance condition (3) is satisfied for every node. In this paper we propose two algorithms on how to do it, both iterative: of the Newton type and a relaxation method. After node position optimization the grid is called invariant.

In this study we consider a class of dissipative systems, i.e. such systems for which there exists a global convex Lyapunov function $G$ (thermodynamic potential) which implements the second law of thermodynamics. For example, for this reason, all reactions on figure 1 are reversible. Dissipative systems have the only steady state in the equilibrium point and as
the time $t$ tends to infinity, the system reaches the equilibrium state, while in the course of the transition the Lyapunov function decreases monotonically. Thermodynamic properties of dissipative systems help a lot: for example, they unambiguously define metrics in the phase space to perform geometrical calculations and also define the choice of projector $P$ almost uniquely (see next section).

Low-dimensional invariant manifolds also exist for systems with a more complicated dynamic behaviour so why study the invariant manifolds of slow motions for a particular class of purely dissipative systems? The answer is the following: most of the physically significant models include non-dissipative components in the form of either a conservative dynamics or in the form of external fluxes. For example, one can think of irreversible reactions among the suggested stoichiometric mechanism (inverse processes are so improbable that we discard them completely, thereby effectively'opening' the system to the remaining irreversible flux). For all such systems, the MIG is applicable almost without special refinements, and bears the significance that invariant manifolds are constructed as a 'deformation' of the relevant manifolds of slow motion of the purely dissipative dynamics. An example of this construction for open systems is presented below in the last section of the paper. The calculations in the last section do not use grid specifics and can be applied not only for grid representation of the invariant manifold, but also for any analytical form of its representation.

## 2 Dissipative Systems and Thermodynamic Projector

### 2.1 Kinetic Equations

Let us introduce the notions used in the paper [see also 3,7,9]. We will consider a closed system with $n$ chemical species $\mathrm{A}_{1}, \ldots, \mathrm{~A}_{n}$, participating in a complex reaction. The complex reaction is represented
of the system. For instance, it becomes immediately apparent in the hydrogen burning model that, for instance, the concentration of molecular hydrogen $\left(\mathrm{H}_{2}\right)$ on the way to equilibrium changes in two stages: very little during the initial fast stage but gradually as equilibrium is approached. In contrast, the oxygen radical ( 0 ) coordinate is fast. Concentration changes rapidly in the initial stage but once equilibrium is approached it reaches a plateau. The behaviour of the hydroxy radical, OH , is deceptive. Depending on the initial conditions it can behave as slow or as fast variable and this fact is clearly visible on the map. The extension of the MIG approach to more sophisticated, yet related reactions in biochemistry or atmospheric chemistry could provide new insights into the behaviour of crucial chemical species.

The researchers conclude that the great utility of their technique for chemical systems may also be applicable beyond chemical kinetics. Their discrete invariant object could be used to visualize other types of functions and properties, for example, in the 'kinetics + transport' systems or in the control theory.

David Bradley of Sciencebase.com
by the following stoichiometric mechanism:
$\alpha_{s 1} \mathrm{~A}_{1}+\ldots+\alpha_{s n} \mathrm{~A}_{n} \rightleftharpoons \beta_{s 1} \mathrm{~A}_{1}+\ldots+\beta_{s n} \mathrm{~A}_{n}$,
where the index $s=1, \ldots, r$ enumerates the reaction steps, and where integers, $\alpha_{s i}$ and $\beta_{s i}$, are stoichiometric coefficients. For each reaction step $s$, we introduce $n$-component vectors $\boldsymbol{\alpha}_{s}$ and $\boldsymbol{\beta}_{s}$ with components $\alpha_{s i}$ and $\beta_{s i}$. Notation $\boldsymbol{\gamma}_{S}$ stands for the vector with integer components $\gamma_{s i}=\beta_{s i}-\alpha_{s i}$ (the stoichiometric vector).

For every $A_{i}$ an extensive variable $N_{i}$, 'the number of particles of that species', is defined. The concentration of $A_{i}$ is $x_{i}=$ $N_{i} / V$, where $V$ is the volume.

Given the stoichiometric mechanism (4), the reaction kinetic equations read:

$$
\begin{equation*}
\dot{\boldsymbol{N}}=V \boldsymbol{J}(\boldsymbol{x}), \boldsymbol{J}(\boldsymbol{x})=\sum_{s=1}^{r} \gamma_{s} W_{s}(\boldsymbol{x}) \tag{5}
\end{equation*}
$$

where the dot denotes the time derivative, and $W_{s}$ is the reaction rate function of the step $s$. In particular, the mass action law suggests the polynomial form of the reaction rates:

$$
\begin{align*}
& W_{s}(\boldsymbol{x})=W_{s}^{+}(\boldsymbol{x})-W_{s}^{-}(\boldsymbol{x})= \\
& k_{s}^{+}(T) \prod_{i=1}^{n} x_{i}^{\alpha_{i}}-k_{s}^{-}(T) \prod_{i=1}^{n} x_{i}^{\beta_{i}} \tag{6}
\end{align*}
$$

where $k_{s}^{+}(T)$ and $k_{s}^{-}(T)$ are the constants of the direct and of the inverse reactions rates of the $s$-th reaction step, $T$ is the temperature.

The rate constants are not independent. The principle of detail balance gives the following connection between these constants: there exists such a positive vector $x^{\mathrm{eq}}(T)$ that

$$
\begin{equation*}
W_{s}^{+}\left(\boldsymbol{x}^{\mathrm{eq}}\right)=W_{s}^{-}\left(\boldsymbol{x}^{\mathrm{eq}}\right) \text { for all } s=1, \ldots, r \tag{7}
\end{equation*}
$$

For $V, T=$ const we do not need additional equations and data. It is possible simply to divide equation (5) by the constant volume and to write

$$
\begin{equation*}
\dot{\boldsymbol{x}}=\sum_{s=1}^{r} \gamma_{s} W_{s}(\boldsymbol{x}) . \tag{8}
\end{equation*}
$$

Conservation laws (balances) impose linear constraints on admissible vectors $\boldsymbol{x}$ :

$$
\begin{equation*}
\left(\boldsymbol{b}_{i}, \boldsymbol{x}\right)=B_{i}=\text { const }, i=1, \ldots, l, \tag{9}
\end{equation*}
$$

where $\boldsymbol{b}_{i}$ are fixed and linearly independent vectors. Let us denote as $\boldsymbol{B}$ the set of vectors which satisfy the conservation laws (9) with given $B_{i}$ :

$$
\boldsymbol{B}=\left\{\boldsymbol{x} \mid\left(\boldsymbol{b}_{1}, \boldsymbol{x}\right)=B_{1}, \ldots,\left(\boldsymbol{b}_{l}, \boldsymbol{x}\right)=B_{l /}\right\} .
$$

The natural phase space $\boldsymbol{X}$ of the system (8) is the intersection of the cone of $n$-dimensional vectors with non-negative components, with the set $\boldsymbol{B}$, and $\operatorname{dim} \boldsymbol{X}=d=$ $n-l$. In addition, we assume that each of the conservation laws is supported by each elementary reaction step, that is

$$
\begin{equation*}
\left(\boldsymbol{\gamma}_{s}, \boldsymbol{b}_{i}\right)=0 \tag{10}
\end{equation*}
$$

for each pair of vectors $\boldsymbol{\gamma}_{S}$ and $\boldsymbol{b}_{i}$.
We assume that the kinetic equation (8) describes evolution towards the unique equilibrium state, $\boldsymbol{x}^{\mathrm{eq}}$, in the interior of the phase space $\boldsymbol{X}$. Furthermore, we assume that there exists a strictly convex function $G(\boldsymbol{x})$ which decreases monotonically in time due to (8):

$$
\begin{equation*}
\dot{G}=(\nabla G(x), J(x)) \leq 0 \tag{11}
\end{equation*}
$$

Here $\nabla G$ is the vector of partial derivatives $\partial G / \partial x_{i}$, and the convexity assumes that the $n \times n$ matrices

$$
\begin{equation*}
\boldsymbol{H}_{\boldsymbol{x}}=\left\|\boldsymbol{\partial}^{2} G(\boldsymbol{x}) / \boldsymbol{\partial} x_{i} \boldsymbol{\partial} x_{j}\right\| \tag{12}
\end{equation*}
$$

are positive definite for all $x \in X$. In addition, we assume that the matrices (12) are invertible if $x$ is taken in the interior of the phase space.

The matrix $\boldsymbol{H}$ defines an important Riemann structure on the concentration space, the thermodynamic (or entropic) scalar product:

$$
\begin{equation*}
\langle\boldsymbol{x}, \boldsymbol{y}\rangle_{c}=\left(\boldsymbol{x}, \boldsymbol{H}_{x} \boldsymbol{y}\right) \tag{13}
\end{equation*}
$$

This choice of the Riemann structure is unambiguous from the thermodynamic perspective. We use this metrics for all geometrical constructions, for measuring angles and distances in the phase space $U$.

The function $G$ is the Lyapunov function of the system (5), and $x^{e q}$ is the point of global minimum of the function $G$ in the phase space $\boldsymbol{X}$. Stated differently, the manifold of equilibrium states $\boldsymbol{x}^{\mathrm{eq}}\left(B_{1}, \ldots, B_{l}\right)$ is the solution to the variational problem,

$$
\begin{equation*}
G \rightarrow \min \text { for }\left(\boldsymbol{b}_{i}, \boldsymbol{x}\right)=B_{i}, \mathrm{i}=1, \ldots, l . \tag{14}
\end{equation*}
$$

For each fixed value of the conserved quantities $B_{i}$, the solution is unique.

For perfect systems in a constant volume under a constant temperature, the Lyapunov function $G$ reads:

$$
\begin{equation*}
G=\sum_{i=1}^{n} x_{i}\left[\ln \left(x_{i} / x_{i}^{\mathrm{eq}}\right)-1\right] . \tag{15}
\end{equation*}
$$

### 2.2 Thermodynamic Projector

For dissipative systems, we keep in mind the following picture (fig.2). The vector field $J(x)$ generates the motion on the phase space $U: d x / d t=J(x)$. An ansatz manifold $\Omega$ is given, it is the current approximation to the invariant manifold. This manifold $\Omega$ is described as the image of the map $F: W \rightarrow U$, where $W$ is a space of macroscopic variables, and $U$ is our phase space.

The projected vector field $P J(x)$ belongs to the tangent space $T_{x}$, and the equation $d x / d t=P J(x)$ describes the motion along the ansatz manifold $\Omega$ (if the initial state belongs to $\Omega$ ). The induced dynamics on the space $W$ is generated by the vector field

$$
\frac{d y}{d t}=\left(D_{y} F\right)^{-1} P J(F(y))
$$

Here the inverse linear operator $\left(D_{y} F\right)^{-1}$ is defined on the tangent space $T_{F(y)}$, because the map $F$ is assumed to be immersion, that is the differential $\left(D_{y} F\right)$ is the isomorphism onto the tangent space $T_{F(y)}$.

Projection operators $P$ contribute to the invariance equation (3). Limiting the results, exact solutions only weakly depend on the particular choice of projectors, or do not depend on it at all. However, thermodynamic validity of approximations obtained at each iteration step towards the limit strongly depends on the choice of the projector.

Let some (not obligatory invariant) manifold $\Omega$ be considered as a manifold of reduced description. We should define a field of linear operators, $\boldsymbol{P}_{x}$, labelled by the states $\boldsymbol{x} \in \Omega$, which project the vectors $J(x), \boldsymbol{x} \in \Omega$ onto the tangent bundle of the manifold $\Omega$, thereby generating the induced vector field, $\boldsymbol{P}_{\boldsymbol{x}} \boldsymbol{J}(\boldsymbol{x}), \boldsymbol{x} \in \Omega$. This induced vector field on the tangent bundle of the manifold $\Omega$ is identified with the reduced dynamics along the manifold $\Omega$. The thermodynamicity requirement for this induced vector field reads

$$
\begin{equation*}
\left(\nabla G(\boldsymbol{x}), \boldsymbol{P}_{\boldsymbol{x}} J(\boldsymbol{x})\right) \leq 0, \text { for each } x \in \Omega \tag{16}
\end{equation*}
$$

The condition (16) means that the entropy $S$ (which is the Lyapunov function with a minus sign) should increase in the new dynamics (2).

How to construct the projector $P$ ? Another form of this question is: how to define the plain of fast motions $x+\operatorname{ker} P$ ? The choice of the projector $P$ is ambiguous, from the formal point of view, but the second law of thermodynamics gives a good idea [3]: the entropy should grow in the fast motion, and the point $x$ should be the point of entropy maximum on the plane of fast motion $x+\operatorname{ker} P$. That is, the subspace ker $P$ should belong to the kernel of the entropy differential:

```
ker }\mp@subsup{P}{x}{}\subset\operatorname{ker}\mp@subsup{D}{x}{}S
```

Of course, this rule is valid for closed systems with entropy, but it can also be extended onto open systems: the projection of the 'thermodynamic part' of $J(x)$ onto $T_{x}$ should have a positive entropy production. If this thermodynamic requirement is valid for any ansatz manifold not tangent to
the entropy levels and for any thermodynamic vector field, then the thermodynamic projector is unique [10]. Let us describe this projector $P$ for a given point $x$, subspace $T_{x}=\mathrm{imP}$, differential $D_{x} S$ of the entropy $S$ at the point $x$ and the second differential of the entropy at the point $x$, the bilinear functional $\left(D_{x}^{2} S\right)_{x}$. We need the positively definite bilinear form $\langle z \mid p\rangle_{x}=$ $-\left(D_{x}^{2} S\right)_{x}(z, p)$ (the entropic scalar product). There exists a unique vector $g$ such that $\langle g \mid p\rangle_{x}=D_{x} S(p)$. It is the Riesz representation of the linear functional $D_{x} S$ with respect to entropic scalar product. If $g \neq 0$ then the thermodynamic projector is

$$
\begin{equation*}
P(J)=P^{\perp}(J)+\frac{g^{\|}}{\langle g \| \mid g\rangle_{x}}\left\langle g^{\perp} \mid J\right\rangle_{x} \tag{17}
\end{equation*}
$$

where $P^{\perp}$ is the orthogonal projector onto $T_{x}$ with respect to the entropic scalar product, and the vector $g$ is split into tangent and orthogonal components:

$$
g=g^{\|}+g^{\perp} ; g^{\|}=P^{\perp} g ; g^{\perp}=\left(1-P^{\perp}\right) g
$$

This projector is defined if $g^{\|} \neq 0$.
If $g=0$ (the equilibrium point) then $P(J)=P^{\perp}(J)$.

For given $T_{x}$, the thermodynamic projector (17) depends on the point $x$ through the $x$-dependence of the scalar product $\langle\mid\rangle_{x}$, and also through the differential of $S$ in $x$.

### 2.3 Symmetric Linearization

The invariance condition (3) supports a lot of invariant manifolds, and not all of them are relevant to the reduced description (for example, any individual trajectory is itself an invariant manifold). This should be carefully taken into account when deriving a relevant equation for the correction in the states of the initial manifold $\Omega_{0}$ which are located far from equilibrium. This point concerns the procedure of the linearization of the vector field $\boldsymbol{J}$, appearing in equation (1). Let $\boldsymbol{c}$ be an arbitrary point of the phase space. The linearization of the vector function $\boldsymbol{J}$ about $\boldsymbol{c}$ may be written $J(\boldsymbol{c}+\delta \boldsymbol{c}) \approx \boldsymbol{J}(\boldsymbol{c})+\boldsymbol{L}_{c} \delta \boldsymbol{c}$ where
the linear operator $\boldsymbol{L}_{\boldsymbol{c}}$ acts as follows (for the mass action law):

$$
\begin{align*}
\boldsymbol{L}_{\boldsymbol{c}} \boldsymbol{x}= & \sum_{s=1}^{r} \boldsymbol{\gamma}_{s}\left[W_{s}^{+}(\boldsymbol{c})\left(\boldsymbol{\alpha}_{s}, \boldsymbol{H}_{\boldsymbol{c}} \boldsymbol{x}\right)\right. \\
& \left.-W_{s}^{-}(\boldsymbol{c})\left(\boldsymbol{\beta}_{s}, \boldsymbol{H}_{\boldsymbol{c}} \boldsymbol{x}\right)\right] \tag{18}
\end{align*}
$$

Here $\boldsymbol{H}_{\boldsymbol{c}}$ is the matrix of second derivatives of the function $G$ in the state $\boldsymbol{c}$, see (12). The matrix $\boldsymbol{L}_{\boldsymbol{c}}$ in (18) can be decomposed as follows:

$$
\begin{equation*}
L_{c}=L_{c}^{\prime}+L_{c}^{\prime \prime} \tag{19}
\end{equation*}
$$

Matrices $\boldsymbol{L}_{c}^{\prime}$ and $\boldsymbol{L}_{c}^{\prime \prime}$ act as follows:
$\boldsymbol{L}_{\boldsymbol{C}}^{\prime} \boldsymbol{x}=-\frac{1}{2} \sum_{s=1}^{r}\left[W_{s}^{+}(\boldsymbol{c})+W_{s}^{-}(\boldsymbol{c})\right] \gamma_{s}\left(\gamma_{s}, \boldsymbol{H}_{\boldsymbol{C}} \boldsymbol{x}\right)$,
$\boldsymbol{L}_{\boldsymbol{C}}^{\prime \prime} \boldsymbol{x}=\frac{1}{2} \sum_{s=1}^{r}\left[W_{s}^{+}(\boldsymbol{c})-W_{s}^{-}(\boldsymbol{c})\right] \gamma_{s}\left(\boldsymbol{\alpha}_{s}+\boldsymbol{\beta}_{s}, \boldsymbol{H}_{\boldsymbol{C}} \boldsymbol{x}\right)$.
$(20,21)$
Some features of this decomposition are best seen when we use the thermodynamic scalar product (13): The following properties of the matrix $\boldsymbol{L}_{\boldsymbol{c}}^{\prime}$ are verified immediately:
(i) The matrix $\boldsymbol{L}_{\boldsymbol{c}}^{\prime}$ is symmetric in the scalar product (13):

$$
\begin{equation*}
\left\langle\boldsymbol{x}, \boldsymbol{L}_{\boldsymbol{c}}^{\prime} \boldsymbol{y}\right\rangle=\left\langle\boldsymbol{y}, \boldsymbol{L}_{\boldsymbol{c}}^{\prime} \boldsymbol{x}\right\rangle \tag{22}
\end{equation*}
$$

(ii) The matrix $\boldsymbol{L}_{\boldsymbol{c}}^{\prime}$ is non-positive definite in the scalar product (13):

$$
\begin{equation*}
\left\langle x, L_{c}^{\prime} x\right\rangle \leq 0 \tag{23}
\end{equation*}
$$

(iii) The null space of the matrix $\boldsymbol{L}_{\boldsymbol{c}}^{\prime}$ is the linear envelope of the vectors $\boldsymbol{H}_{c}^{-1} \boldsymbol{b}_{i}$ representing the complete system of conservation laws:

$$
\begin{equation*}
\operatorname{ker} \boldsymbol{L}_{\boldsymbol{c}}^{\prime}=\operatorname{Lin}\left\{\boldsymbol{H}_{\boldsymbol{c}}^{-1} \boldsymbol{b}_{i}, i=1, \ldots, l\right\} \tag{24}
\end{equation*}
$$

(iv) If $c=c^{\text {eq }}$, then

$$
W_{s}^{+}\left(\boldsymbol{c}^{\mathrm{eq}}\right)=W_{s}^{-}\left(\boldsymbol{c}^{\mathrm{eq}}\right), \text { and }
$$

$$
\begin{equation*}
\boldsymbol{L}_{\boldsymbol{C}^{\mathrm{eq}}}^{\prime}=\boldsymbol{L}_{\boldsymbol{C}^{\mathrm{eq}}} \tag{25}
\end{equation*}
$$

Thus, the decomposition (19) splits the matrix $L_{c}$ in two parts: one part, (20), is symmetric and non-positive definite, while the other part, (21), vanishes in the equi-
librium. The decomposition (19) explicitly takes into account the mass action law. For other dissipative systems, the decomposition (19) is possible as soon as the relevant kinetic operator is written in a gain-loss form.

## 3 Invariant Grids

In most of the works (of us and of other people on similar problems), analytic forms were required to represent manifolds (see, however, the method of Legendre integrators [11-13]). However, in order to construct manifolds of a relatively low dimension, grid-based representations of manifolds become a relevant option [8].

The main idea of the MIG is to find a mapping of the finite-dimensional grids into the phase space of a dynamic system. That is, we construct not just a point approximation of the invariant manifold $F^{*}(y)$, but an invariant grid. When refined, it is expected to converge, of course, to $F^{*}(y)$, but in any case it is a separate, independently defined object.

Let us denote $L=R^{n}, G$ is a discrete subset of $R^{n}$. It is natural to think of a regular grid, but this is not so crucial. For every point $y \in G$, a neighbourhood of $y$ is defined: $V_{y} \subset G$, where $V_{y}$ is a finite set, and, in particular, $y \in V_{y}$. On regular grids, $V_{y}$ includes, as a rule, the nearest neighbours of $y$. It may also include the points next to the nearest neighbours.

For our purpose, we need to define a grid differential operator. For every function, defined on the grid, all derivatives are also defined:

$$
\begin{equation*}
\left.\frac{\partial f}{\partial y_{i}}\right|_{y \in G}=\sum_{z \in V_{y}} q_{i}(z, y) f(z), i=1, \ldots n, \tag{26}
\end{equation*}
$$

where $q_{i}(z, y)$ are some coefficients.
Here we do not specify the choice of the functions $q_{i}(z, y)$. We just mention in passing that, as a rule, equation (26) is established using some approximation of $f$ in the neighbourhood of $y$ in $R^{n}$ by some differentiable functions (for example, poly-
nomials). This approximation is based on the values of $f$ at the points of $V_{y}$. For regular grids, $q_{i}(z, y)$ are functions of the difference $z-y$. For some of the nodes $y$ which are close to the edges of the grid, functions are defined only on the part of $V_{y}$. In this case, the coefficients in (26) should be modified appropriately in order to provide an approximation using the available values of $f$. Below we assume this modification is always done. We also assume that the number of points in the neighbourhood $V_{y}$ is always sufficient to make the approximation possible. This assumption restricts the choice of the grids G. Let us call admissible all such subsets $G$, on which one can define differentiation operator in every point.

Let $F$ be a given mapping of some admissible subset $G \subset R^{n}$ into $U$. For every $y \in V$ we define tangent vectors:

$$
\begin{equation*}
T_{y}=\operatorname{Lin}\left\{g_{i}\right\}_{1}^{n}, \tag{27}
\end{equation*}
$$

where vectors $g_{i}(i=1, \ldots, n)$ are partial derivatives (26) of the vector function $F$ :

$$
\begin{equation*}
g_{i}=\frac{\partial F}{\partial y_{i}}=\sum_{z \in V_{y}} q_{i}(z, y) F(z), \tag{28}
\end{equation*}
$$

or in the coordinate form:

$$
\begin{equation*}
\left(g_{i}\right)_{j}=\frac{\partial F_{j}}{\partial y_{i}}=\sum_{z \in V_{y}} q_{i}(z, y) F_{j}(z) . \tag{29}
\end{equation*}
$$

Here $\left(g_{i}\right)_{j}$ is the $j$-th coordinate of the vector $\left(g_{i}\right)$, and $F_{j}(z)$ is the $j$-th coordinate of the point $F(z)$.

The grid $G$ is invariant, if for every node $y \in G$ the vector field $J(F(y))$ belongs to the tangent space $T_{y}$ (here $J$ is the right hand side of the kinetic equation (1)).

So, the definition of the invariant grid includes:

1) The finite admissible subset $G \subset R^{n}$.
2) A mapping $F$ of this admissible subset $G$ into $U$ (where $U$ is the phase space of kinetic equation (1)).
3) The differentiation formulas (26) with given coefficients $q_{i}(z, y)$.

The grid invariance equation has the form of an inclusion:

## $J(F(y)) \in T_{y}$ for every $y \in G$,

or the form of an equation:

$$
\left(1-P_{y}\right) J(F(y))=0 \text { for every } y \in G
$$

where $P_{y}$ is the thermodynamic projector (17).

The grid differentiation formula (26) is needed, in the first place, to establish the tangent space $T_{y}$, and the null space of the thermodynamic projector $P_{y}$ in each node. It is important to realize that the locality of the construction of the thermodynamic projector makes this possible without a global parametrization.
Let $x=F(y)$ be the location of the grid's node $y$ immersed into $U$. We have the set of tangent vectors $g_{i}(x)$, defined in $x(28)$, (29). Thus, the tangent space $T_{y}$ is defined by (27). Also, we have the entropy function $S(x)$, the linear functional $D_{x} S S_{x}$, and the subspace $T_{0 y}=\left.T_{y} \cap \operatorname{ker} D_{x} S\right|_{x}$ in $T_{y}$. Let $T_{0 y} \neq T_{y}$. In this case we have a vector $\boldsymbol{e}_{y} \in$ $T_{y}$, orthogonal to $T_{0 y}, D_{x} S S_{x}\left(\boldsymbol{e}_{y}\right)=1$.

Then the thermodynamic projector is defined as:

$$
\begin{equation*}
P_{y} \cdot=P_{0 y} \cdot+e_{y} D_{x} S_{x} \cdot \tag{30}
\end{equation*}
$$

where $P_{0 y}$ is the orthogonal projector on $T_{0 y}$ with respect to the entropic scalar product $\langle\mid\rangle_{x}$.

If $T_{0 y}=T_{y}$, then the thermodynamic projector is the orthogonal projector on $T_{y}$ with respect to the entropic scalar product $\langle\mid\rangle_{x}$.

The general schema of solving the invariance equation (3) to optimize positions of the invariant grid nodes in space is the following:
0) The grid is initialized. For example, it is possible to use spectral decomposition of $\left(D_{x}^{2} S\right) x$ in the equilibrium.

1) Given some node positions, the tangent vectors in every node of the grid (27) are calculated; at this stage the connectivity between nodes is used.
2) With set of tangent vectors calculated at the previous step, solve the invariance equation for every node independently and
calculate a shift $\delta y$ of every node in the phase space; we propose two algorithms to calculate the shift: the Newton method with incomplete linearization and the relaxation method [see also 4-6,8].
3) Repeat steps 1 and 2 until some convergence criterion will be fulfilled: for example, all shifts $\delta y_{i}, i=1, \ldots, n$ will be less than a predefined $\varepsilon_{\text {conv }}$.
4) Update the structure of the grid: for example, add new nodes and extend (extrapolate) or refine (interpolate) the grid. Some strategies for this are described further.
5) Repeat steps $1-4$ until some criterion will be fulfilled: typically, when the nodes reach the phase space boundary or the spectral gap is too small (see below).

The idea of the Newton method with incomplete linearization is to use linear approximation of $J$ in the vicinity of a grid node $y$ (keeping the projector $P$ fixed). At the same time the node is shifted in the fast direction (in $y+$ ker $P_{y}$ affine subspace).

For the Newton method with incomplete linearization, the equations for calculation of the new node location $y^{\prime}=y+\delta y$ are:

$$
\left\{\begin{array}{l}
P_{y} \delta y=0  \tag{31}\\
\left(1-P_{y}\right)(J(y)+D J(y) \delta y)=0 .
\end{array}\right.
$$

Here $D J(y)$ is a matrix of derivatives of $J$ evaluated at $y$. Instead of $D J(y)$ (especially in the regions that are far from the equilibrium) the symmetric operator $\boldsymbol{L}^{\prime}(y)$ (20) can be used; this will provide better convergence towards the 'true' invariant manifold.

Equation (31) is a system of linear algebraic equations. In practice, it proves convenient to choose some orthonormal (with respect to the entropic scalar product) basis $\boldsymbol{b}_{i}$ in ker $P_{y}$. Let $r=\operatorname{dim}\left(\operatorname{ker} P_{y}\right)$. Then $\delta y=\sum_{i=1}^{r} \delta_{i} \boldsymbol{b}_{i}$, and system (31) takes the form
$\sum_{k=1}^{r} \delta_{k}\left\langle\boldsymbol{b}_{i} \mid D J(y) \boldsymbol{b}_{k}\right\rangle_{y}=-\left\langle J(y) \mid \boldsymbol{b}_{i}\right\rangle_{y}, i=1 \ldots r$.

This is the system of linear equations for adjusting the node location according to the Newton method with incomplete linearization. We stress once again that one should use the entropic scalar products.

For the relaxation method, one needs to calculate the defect $\Delta_{y}=\left(1-P_{y}\right) J(y)$, and the relaxation step

$$
\begin{equation*}
\tau(y)=-\frac{\left\langle\Delta_{y} \mid \Delta_{y}\right\rangle_{y}}{\left\langle\Delta_{y} \mid D J(y) \Delta_{y}\right\rangle_{y}} \tag{33}
\end{equation*}
$$

Then, the new node location $y^{\prime}$ is computed as

$$
\begin{equation*}
y^{\prime}=y+\tau(y) \Delta_{y} . \tag{34}
\end{equation*}
$$

This is the equation for adjusting the node location according to the relaxation method.

## 4 Grid Construction Strategy

Of all the reasonable strategies of the invariant grid construction we consider here the following two: the growing lump and the invariant flag.

### 4.1 Growing Lump

The construction is initialized from the equilibrium point $y^{*}$. The first approximation is constructed as $F\left(y^{*}\right)=x^{*}$, and for some initial $V_{0}\left(V_{y^{*}} \subset V_{0}\right)$ one has $F(y)=$ $x^{*}+A\left(y-y^{*}\right)$, where $A$ is an isometric embedding (in the standard Euclidean metrics) of $R^{n}$ in $E$.

For this initial grid one makes a fixed number of iterations of one of the methods chosen (Newton's method with incomplete linearization or the relaxation method), and, after that, puts $V_{1}=\cup_{y} \in_{V_{0}} V_{y}$ and extends $F$ from $V_{0}$ onto $V_{1}$ using the linear extrapolation, and the process continues. One of the possible variants of this procedure is to extend the grid from $V_{i}$ to $V_{i+1}$ not after a fixed number of iterations, but only after the invariance defect $\Delta_{y}$ becomes less than a given $\varepsilon$ (in a given norm, which is entropic, as a rule), for all nodes $y \in V_{i}$. The lump stops growing after it reaches the boundary and is within a given accuracy $\|\Delta\|<\varepsilon$.

### 4.2 Invariant Flag

In order to construct the invariant flag sufficiently regular grids $G$ are used, in which many points are located on the coordinate lines, planes, etc. The standard flag $R^{0} \subset R^{1} \subset R^{2} \subset \ldots \subset R^{n}$ (every next space is constructed by adding one more coordinate) is considered. It corresponds to a sequence of grids $\left\{y^{*}\right\} \subset G^{1} \subset G^{2} \ldots \subset$ $G^{n}$, where $\left\{y^{*}\right\}=R^{0}$, and $G^{i}$ is a grid in $R^{i}$.

First, $y^{*}$ is mapped on $x^{*}$ and further $F\left(y^{*}\right)=x^{*}$. Then the invariant grid is constructed on $V^{1} \subset G^{1}$ (up to the boundaries and within a given accuracy $\|\Delta\|<\varepsilon$ ). After that, the neighbourhoods in $G^{2}$ are added to the points $V^{1}$, and the grid $V^{2} \subset G^{2}$ is constructed (up to the boundaries and within a given accuracy) and so on, until $V^{n} \subset G^{n}$ is constructed.

While constructing the $k$-th-order grid $V^{k} \subset G^{k}$, the important role of the grids of lower dimension $V^{0} \subset \ldots \subset V^{k-1} \subset V^{k}$ embedded in it is preserved. The point $F\left(y^{*}\right)=x^{*}$ (equilibrium) remains fixed. For every $y \in V^{q}(q<k)$ the tangent vectors $g_{1}, \ldots, g_{q}$ are constructed, using the differentiation operators (26) on the whole $V^{k}$. Using the tangent space $T_{y}=$ $\operatorname{Lin}\left\{g_{1}, \ldots, g_{q}\right\}$, the projector $P_{y}$ is constructed, the iterations are applied and so on. All this is done in order to obtain a sequence of embedded invariant grids, given by the same map $F$.

### 4.3 Boundaries Check and Entropy

We construct grid mapping of $F$ onto a finite set $V \in G$. The technique of checking whether the grid still belongs to the phase space $U$ of the kinetic system $(F(V) \subset U)$ is quite straightforward: all the points $y \in$ $V$ are checked whether they belong to $U$. If at the next iteration a point $F(y)$ leaves $U$, then it is pulled inside by a homothety transform with the center in $x^{*}$. Since the entropy is a concave function, the homothety contraction with the center in $x^{*}$ increases the entropy monotonically. Another variant to cut off the points which leave $U$.

By construction (17), the kernel of the entropic projector is annulled by the entropy differential. Thus, in the first order, the steps in the Newton method with incomplete linearization (31) as well as in the relaxation method (33) do not change the entropy. But if the steps are quite large, then the increase of the entropy may become essential, and the points are returned on their entropy levels by the homothety contraction with the centre in the equilibrium point.

## 5 Instability of Fine Grids

When one reduces the grid spacing in order to refine the grid, then, once the grid spacing becomes small enough, one can face the problem of the Courant instability [14-16].Instead of converging, at every iteration the grid becomes more and more entangled (see fig. 3).

A way to avoid such instability is wellknown. This is decreasing the time step. In our problem, instead of a true time step, we have a shift in the Newtonian direction. Formally, we can assign the value $h=1$ for one complete step in the Newtonian direc-


Fig. 3. Grid instability. For small grid steps approximations in the calculation of grid derivatives lead to the grid instability effect. Several successive iterations of the algorithm without adaptation of the time step are shown that lead to undesirable 'oscillations', which eventually destroy the grid starting from one of its ends.
tion. Let us extend now the Newton method to arbitrary $h$. For this, let us find $\delta x=$ $\delta F(y)$ from (31), but update $\delta x$ proportionally to $h$; the new value of $x_{n+1}=$ $F_{n+1}(y)$ is equal to

$$
\begin{equation*}
F_{n+1}(y)=F_{n}(y)+h_{n} \delta F_{n}(y) \tag{35}
\end{equation*}
$$

where $n$ denotes the number of iteration.
One way to choose the step value $h$ is to make it adaptive, by controlling the average value of the invariance defect $\left\|\Delta_{y}\right\|$ at every step. Another way is the convergence control: then $\Sigma h_{n}$ plays a role of time.

Elimination of the Courant instability for the relaxation method can be done quite analogously. Everywhere the step $h$ is maintained as large as it is possible without running into convergence problems.

## 6 Analyticity and Effect of Superresolution

When constructing invariant grids, one must define the differential operators (26) for every grid's node. For calculating the differential operators in some point $y$, an interpolation procedure in the neighbourhood of $y$ is used. As a rule, it is an interpolation by a low-order polynomial, which is constructed using the function values in the nodes belonging to the neighbourhood of $y$ in $G$. This approximation (using values in the nearest neighbourhood nodes) is natural for smooth functions. But we are looking for the analytical invariant manifold. Analytical functions have a much more 'rigid' structure than the smooth ones. One can change a smooth function in the neighbourhood of any point in such a way that outside this neighbourhood the function will not change. In general, this is not possible for analytical functions: a kind of a 'long-range' effect takes place (as is well known).

The idea is to make use of this effect and to reconstruct some analytical function $f_{G}$ using a function given on $G$. There is one important requirement: if the values given on $G$ are values of some function $f$ which is analytical in a neighbourhood $U$, then, if
the $G$ is refined 'correctly', one must have $f_{G} \rightarrow f$ in $U$. The sequence of reconstructed function $f_{G}$ should converge to the 'right' function $f$.

What is the 'correct refinement'? For smooth functions for the convergence $f_{G} \rightarrow f$ it is necessary and sufficient that, in the course of refinement, $G$ would approximate the whole $U$ with arbitrary accuracy. For analytical functions it is necessary only that, under the refinement, $G$ would approximate some uniqueness set $A \subset U$. A subset $A \subset U$ is called uniqueness set in $U$ if for analytical functions in $U \psi$ and $\varphi$ from $\left.\left.\psi\right|_{A} \equiv \varphi\right|_{A}$ it follows $\psi \equiv \varphi$. Suppose we have a sequence of grids $G$, each following one is finer than the previous one, which approximates a set $A$. For smooth functions using function values defined on the grids one can reconstruct the function in $A$. For analytical functions, if the analyticity domain $U$ is known, and $A$ is a uniqueness set in $U$, then one can reconstruct the function in $U$. The set $U$ can be essentially bigger than $A$; because of this such an extension was named as superresolution effect [17]. There exist formulas for construction of analytical functions $f_{G}$ for different domains $U$, uniqueness sets $A \subset$ $U$ and for different ways of discrete approximation of $A$ by a sequence of refined grids $G$ [17]. Here we provide only one Carleman's formula which is the most appropriate for our purposes.

Let domain $U=Q_{\sigma}^{n} \subset C^{n}$ be a product of strips $Q_{\sigma} \subset C, Q_{\sigma}=\{z \mid \operatorname{Im} z<\sigma\}$. We shall construct functions holomorphic in $Q_{\sigma}^{n}$. This is effectively equivalent to the construction of real analytical functions $f$ in the whole $R^{n}$ with a condition on the convergence radius $r(x)$ of the Taylor series for $f$ as a function of each coordinate: $r(x) \geq \sigma$ in every point $x \in R^{n}$.

The sequence of refined grids is constructed as follows: let for every $l=1, \ldots, n$ a finite sequence of distinct points $N_{l} \subset Q_{\sigma}$ be defined:

$$
\begin{equation*}
N_{l}=\left\{x_{l j} \mid j=1,2,3 \ldots\right\}, x_{l j} \neq x_{l i} \text { for } i \neq j \tag{36}
\end{equation*}
$$

The countable uniqueness set $A$, which is approximated by a sequence of refined grids, has the form:

$$
\begin{align*}
A & =N_{1} \times N_{2} \times \ldots \times N_{n} \\
& \left.=\left\{\left(x_{1 i_{1}}, x_{2 i_{2}}, \ldots, x_{n i_{n}}\right) i_{1, \ldots, n}=1,2,3, \ldots\right\}\right\} \tag{37}
\end{align*}
$$

The grid $G_{m}$ is defined as the product of initial fragments $N_{l}$ of length $m$ :
$G_{m}=\left\{\left(x_{1 i_{1}}, x_{2 i_{2}} \ldots x_{n i_{n}}\right) \mid 1 \leq i_{1, \ldots, n} \leq m\right\}$
Let us denote $\lambda=2 \sigma / \pi$ ( $\sigma$ is a halfwidth of the strip $Q_{\sigma}$ ). The key role in the construction of the Carleman formula is played by the functional $\omega_{m}^{\lambda}(u, p, l)$ of 3 variables: $u \in U=Q_{\sigma}^{n}, p$ is an integer, $1 \leq p \leq m, l$ is an integer, $1 \leq p \leq n$. Further $u$ will be the coordinate value at the point where the extrapolation is calculated, $l$ will be the coordinate number, and $p$ will be an element of multi-index $\left\{i_{1}, \ldots, i_{n}\right\}$ for the point $\left(x_{1 i_{1}}, x_{2 i_{2}}, \ldots, x_{n i_{\mathrm{n}}}\right) \in G$ :

$$
\begin{align*}
& \omega_{m}^{\lambda}(u, p, l)= \\
& \frac{\left(e^{\lambda x_{l p}}+e^{\lambda \bar{x}_{l p}}\right)\left(e^{\lambda u}-e^{\lambda x_{l p}}\right)}{\lambda\left(e^{\lambda u}+e^{\lambda \bar{x}_{l p}}\right)\left(u-x_{l p}\right) e^{\lambda x_{l p}}} \\
& \times \prod_{j=1 j \neq p}^{m} \frac{\left(e^{\lambda x_{l p}}+e^{\lambda \bar{x}_{l j}}\right)\left(e^{\lambda u}-e^{\lambda x_{l j}}\right)}{\left(e^{\lambda x_{l p}}-e^{\lambda x_{l j}}\right)\left(e^{\lambda u}+e^{\lambda \bar{x}_{l j}}\right)} \tag{39}
\end{align*}
$$

For real-valued $x_{p k}$ formula (39) simplifies:

$$
\begin{array}{r}
\omega_{m}^{\lambda}(u, p, l)=2 \frac{e^{\lambda u}-e^{\lambda x_{l p}}}{\lambda\left(e^{\lambda u}+e^{\lambda x_{l p}}\right)\left(u-x_{l p}\right)} \times \\
\prod_{j=1 j \neq p}^{m} \frac{\left(e^{\lambda x_{l p}}+e^{\lambda x_{l j}}\right)\left(e^{\lambda u}-e^{\lambda x_{l j}}\right)}{\left(e^{\lambda x_{l p}}-e^{\lambda x_{l j}}\right)\left(e^{\lambda u}+e^{\lambda x_{l j}}\right)} \tag{40}
\end{array}
$$

The Carleman formula for extrapolation from $G_{M}$ on $U=Q_{\sigma}^{n}(\sigma=\pi \lambda / 2)$ has the form $\left(z=\left(z_{1}, \ldots, z_{n}\right)\right)$ :

$$
\begin{equation*}
f_{m}(z)=\sum_{k_{1}, \ldots, k_{n}=1}^{m} f\left(x_{k}\right) \prod_{j=1}^{n} \omega_{m}^{\lambda}\left(z_{j}, k_{j}, j\right) \tag{41}
\end{equation*}
$$

where $k=k_{1}, \ldots, k_{n}, x_{k}=\left(x_{1 k_{1}}, x_{2 k_{2}}, \ldots, x_{n k_{\mathrm{n}}}\right)$.
There exists a theorem [17]:

If $f \in H^{2}\left(Q_{\sigma}^{n}\right)$, then $f(z)=$ $\lim _{m \rightarrow \infty} f_{m}(z)$, where $H^{2}\left(Q_{\sigma}^{n}\right)$
is the Hardy class of holomorphic in $Q_{\sigma}^{n}$ functions.

It is useful to present the asymptotics of (41) for large $\left|\operatorname{Re} z_{j}\right|$. For this purpose, we shall consider the asymptotics of (41) for large $|\operatorname{Reu}|$ :
$\left|\omega_{m}^{\lambda}(u, p, l)\right|=\left|\frac{2}{\lambda u} \prod_{j=1 j \neq p}^{m} \frac{e^{\lambda x_{l p}}+e^{\lambda x_{l j}}}{e^{\lambda x_{l p}}-e^{\lambda x_{l j}}}\right|+o\left(|\operatorname{Re} u|^{-1}\right)$.

From the formula (41) one can see that for the finite $m$ and $\left|\operatorname{Re} z_{j}\right| \rightarrow \infty$ function $\left|f_{m}(z)\right|$ behaves like const $\Pi_{j}\left|z_{j}\right|^{-1}$.

This property (zero asymptotics) must be taken into account when using formula (41). When constructing invariant manifolds $F(W)$, it is natural to use (41) not for the immersion $F(y)$, but for the deviation of $F(y)$ from some analytical ansatz $F_{0}(y)$ [18-20].

The analytical ansatz $F_{0}(y)$ can be obtained using Taylor series, just as in the Lyapunov auxiliary theorem [21]. Another variant is to use Taylor series for the construction of Pade approximations.

It is natural to use approximations (41) in terms of dual variables as well, since there exists for them (as the examples demonstrate) a simple and effective linear ansatz for the invariant manifold. This is the slow invariant subspace $E_{\text {slow }}$ of the operator of linearized system (1) in dual variables at the equilibrium point. This invariant subspace corresponds to the set of 'slow' eigenvalues (with small $|\operatorname{Re} \lambda|$, $\operatorname{Re} \lambda<0$ ). In the space of concentrations this invariant subspace is the quasi-equilibrium manifold. It consists of the maximum entropy points on the affine manifolds of the form $x+E_{\text {fast }}$, where $E_{\text {fast }}$ is the 'fast' invariant subspace of the operator of the linearized system (1) at the equilibrium point. It corresponds to the 'fast' eigenvalues (large $|\operatorname{Re} \lambda|, \operatorname{Re} \lambda<0$ ).

Carleman's formulas can be useful for the invariant grid construction in two plac-
es: first, for the definition of the grid differential operators (26), and second, for the analytical continuation of the manifold from the grid.

## 7 Example: Two-Step

## Catalytic Reaction

Let us consider a two-step four-component reaction with one catalyst $A 2$ (the Mi -chaelis-Menten mechanism, see fig. 1a):

$$
\begin{equation*}
A_{1}+A_{2} \leftrightarrow A_{3} \leftrightarrow A_{2}+A_{4} . \tag{43}
\end{equation*}
$$

We assume the Lyapunov function of the form

$$
S=-G=-\sum_{i=1}^{4} c_{i}\left[\ln \left(c_{i} / c_{i}^{\mathrm{eq}}\right)-1\right]
$$

The kinetic equation for the four-component vector of concentrations, $c=\left(c_{1}, c_{2}\right.$, $c_{3}, c_{4}$ ), has the form

$$
\begin{equation*}
\dot{c}=\gamma_{1} W_{1}+\gamma_{2} W_{2} \tag{44}
\end{equation*}
$$

Here $\gamma_{1,2}$ are stoichiometric vectors,

$$
\begin{equation*}
\gamma_{1}=(-1,-1,1,0), \quad \gamma_{2}=(0,1,-1,1), \tag{45}
\end{equation*}
$$

while functions $W_{1,2}$ are reaction rates:
$W_{1}=k_{1}^{+} c_{1} c_{2}-k_{1}^{-} c_{3}, W_{2}=k_{2}^{+} c_{3}-k_{2}^{-} c_{2} c_{4}$.
Here $k_{1,2}^{ \pm}$are reaction rate constants. The system under consideration has two conservation laws,

$$
\begin{equation*}
c_{1}+c_{3}+c_{4}=B_{1}, \quad c_{2}+c_{3}=B_{2} \tag{47}
\end{equation*}
$$

or $\left\langle\boldsymbol{b}_{1,2}, \boldsymbol{c}\right\rangle=B_{1,2}$, where $\boldsymbol{b}_{1}=(1,0,1,1)$ and $\boldsymbol{b}_{2}=(0,1,1,0)$. The non-linear system (43) is effectively 2 -dimensional, and we consider a 1-dimensional reduced description. For our example, we chose the following set of parameters:

$$
\begin{array}{ll}
k_{1}^{+}=0.3, & k_{1}^{-}=0.15, \\
k_{2}^{+}=0.8, & k_{2}^{-}=2.0 ; \\
c_{1}^{\text {eq }}=0.5, & c_{2}^{\text {eq }}=0.1, c_{3}^{\text {eq }}=0.1, c_{4}^{\text {eq }}=0.4 ;  \tag{48}\\
B_{1}=1.0, & B_{2}=0.2
\end{array}
$$

The 1-dimensional invariant grid is shown in figure 4 in the $\left(c_{1}, c_{4}, c_{3}\right)$ coordinates. The grid was constructed by the growing lump method, as described above.


Fig. 4. One-dimensional invariant grid (circles, bold line) for the 2-dimensional chemical system. Projection into the 3 -dimensional space of c1, c4, c3, concentrations. The trajectories of the system in the phase space are shown by lines. The equilibrium point is marked by the square. The system quickly reaches the grid and further moves along it.

We used Newton iterations to adjust the nodes. The grid was grown up to the boundaries of the phase space.

The grid in this example is a 1-dimensional ordered sequence $\left\{x_{1}, \ldots, x_{n}\right\}$. The grid derivatives for calculating the tangent vectors $g$ were taken as $g\left(x_{i}\right)=$ $\left(x_{i+1}-x_{i-1}\right) /\left\|x_{i+1}-x_{i-1}\right\|$ for the internal nodes, and $g\left(x_{1}\right)=\left(x_{1}-x_{2}\right) /\left\|x_{1}-x_{2}\right\|$, $g\left(x_{n}\right)=\left(x_{n}-x_{\mathrm{n}-1}\right) /\left\|x_{n}-x_{n-1}\right\|$ for the grid's boundaries.

Close to the phase space boundaries we had to apply an adaptive algorithm for choosing the time step $h$ : if, after the next growing step (adding new nodes to the grid and after completing $N=20$ Newtonian steps), the grid did not converge, then we chose a new step size $h_{n+1}=h_{n} / 2$ and recalculate the grid. The final (minimal) value for $h$ was $h \approx 0.001$.

The location of the nodes was parametrized with the entropic distance to the equilibrium point measured in the quadratic metrics given by the matrix $\boldsymbol{H}_{c}=$ $-\left\|\partial^{2} S(c) / \partial c_{i} \partial c_{j}\right\|$ in the equilibrium $c^{\text {eq }}$. It means that every node is located on a sphere in this metrics with a given radius, which increases linearly with the number of the node. In figure 4 the step of the in-
crease is chosen to be 0.05 . Thus, the first node is at the distance of 0.05 from the equilibrium, the second is at the distance of 0.10 and so on. Figure 5 shows several important quantities which facilitate the understanding of the object (invariant grid) extracted. The sign on the x-axis of the graphs in figure 5 is meaningless since the distance is always positive, but in this
situation it indicates two possible directions from the equilibrium point.

Figure 5 a , b represents the slow 1-dimensional component of the dynamics of the system. Given any initial condition, the system quickly finds the corresponding point on the manifold and starting from this point the dynamics is given by part of the graph in figure $5 \mathrm{a}, \mathrm{b}$.


Fig. 5. One-dimensional invariant grid for the 2-dimensional chemical system. a Values of the concentrations along the grid. b Values of the entropy and the entropy production $(-d G / d t)$ along the grid. c Ratio of the relaxation times 'towards' and 'along' the manifold. The node positions are parametrized with entropic distance measured in the quadratic metrics iven by $\boldsymbol{H}_{c}=-\left\|\partial^{2} S(c) / \partial c_{i} \partial c_{j}\right\|$ in the equilibrium $c^{\mathrm{eq}}$. Entropic coordinate equal to zero corresponds to the equilibrium.

One of the useful quantities is shown on the figure 5 c . It is the relation between the relaxation times 'toward' and 'along' the $\operatorname{grid}\left(\lambda_{2} / \lambda_{1}\right.$, where $\lambda_{1}, \lambda_{2}$ are the smallest and the next smallest by absolute value non-zero eigenvalue of the system, symmetrically linearized at the point of the grid node). The figure demonstrates that the system is very stiff close to the equilibrium point ( $\lambda_{1}$ and $\lambda_{2}$ are well separated from each other), and becomes less stiff (by order of magnitude) near the boundary. This leads to the conclusion that the 1-dimensional reduced model is more adequate in the neighbourhood of the equilibrium where fast and slow motions are separated by two orders of magnitude. On the end-points of the grid the 1 -dimensional reduction ceases to be well-defined.

## 8 Example: Model Hydrogen Burning Reaction

In this section we consider a more complicated example (see fig. 1b), where the concentration space is 6 -dimensional, while the system is 4 -dimensional. We construct an invariant flag which consists of 1 - and 2-dimensional invariant manifolds.

We consider a chemical system with six species called $\mathrm{H}_{2}$ (hydrogen), $\mathrm{O}_{2}$ (oxygen), $\mathrm{H}_{2} \mathrm{O}$ (water), $\mathrm{H}, \mathrm{O}, \mathrm{OH}$ (radicals) (see fig. 1). We assume the Lyapunov function of the form $S=-G=$

$$
-\sum_{i=1}^{6} c_{i}\left[\ln \left(c_{i} / c_{i}^{\mathrm{eq}}\right)-1\right] .
$$

The subset of the hydrogen burning reaction and corresponding (direct) rate constants were taken as:

| 1. $H_{2} \leftrightarrow 2 H$ | $k_{1}^{+}=2$ |
| :--- | :--- |
| 2. $O_{2} \leftrightarrow 2 O$ | $k_{2}^{+}=1$ |
| 3. $H_{2} O \leftrightarrow H+O H$ | $k_{3}^{+}=1$ |
| 4. $H_{2}+O \leftrightarrow H+O H$ | $k_{4}^{+}=10^{3}$ |
| 5. $O_{2}+H \leftrightarrow O+O H$ | $k_{5}^{+}=10^{3}$ |
| 6. $H_{2}+O \leftrightarrow H_{2} O$ | $k_{6}^{+}=10^{2}$ |

The conservation laws are:
$2 c_{\mathrm{H}_{2}}+2 c_{\mathrm{H}_{2} \mathrm{O}}+c_{\mathrm{H}}+c_{\mathrm{OH}}=b_{\mathrm{H}}$
$2 c_{O_{2}}+c_{\mathrm{H}_{2} \mathrm{O}}+c_{\mathrm{O}}+c_{\mathrm{OH}}=b_{\mathrm{O}}$
For parameter values we took $b_{H}=2$, $b_{O}=1$, and the equilibrium point:

$$
\begin{array}{ll}
c_{c_{2}^{\mathrm{eq}}}^{\mathrm{eq}}=0.27 & c_{O_{2}}^{\mathrm{eq}}=0.135 \\
c_{H_{2} O}^{\mathrm{eq}}=0.7 & c_{H}^{\mathrm{eq}}=0.05 \\
c_{O}^{\mathrm{eq}}=0.02 & c_{O H}^{\mathrm{eq}}=0.01 \tag{51}
\end{array}
$$

Other rate constants $k_{i}^{-}=1 . .6$ were calculated from $c^{\mathrm{eq}}$ value and $k_{i}^{+}$. For this system the stoichiometric vectors are:

$$
\begin{array}{ll}
\gamma_{1}=(-1,0,0,2,0,0) & \gamma_{2}=(0,-1,0,0,2,0) \\
\gamma_{3}=(0,0,-1,1,0,1) & \gamma_{4}=(-1,0,0,1,-1,1) \\
\gamma_{5}=(0,-1,0,-1,1,1) & \gamma_{6}=(-1,0,1,0,-1,0)
\end{array}
$$

The system under consideration is fictitious in the sense that the subset of equa-


Fig. 6. One-dimensional invariant grid for model hydrogen burning reaction. a Projection into the 3 -dimensional space of $c_{\mathrm{H}}, \mathrm{C}_{\mathrm{O}}, \mathrm{COH}_{\mathrm{OH}}$ concentrations. b Concentration values along the grid. c Three smallest by the absolute value non-zero eigenvalues of the symmetrically linearized system.
tions corresponds to the simplified picture of this chemical process and the rate constants do not correspond to any experimentally measured quantities, rather they reflect only orders of magnitudes relevant real-world systems. In that sense we consider here a qualitative model system, which allows us to illustrate the invariant grids method. Nevertheless, modelling of more realistic systems differs only in the
number of species and equations. This leads, of course, to computationally harder problems, but difficulties are not crucial.

Figure 6a presents a 1-dimensional invariant grid constructed for the system. Figure 6b demonstrates the reduced dynamics along the manifold (for the explanation of the meaning of the $x$-coordinate, see the previous subsection). In figure 6 c the three smallest non-zero eigenvalues by


Fig. 7. Two-dimensional invariant grid for the model hydrogen burning reaction. a Projection into the 3 -dimensional space of $c_{\mathrm{H}}, \mathrm{c}_{\mathrm{O}}, \mathrm{C}_{\mathrm{OH}}$ concentrations. b Projection into the principal 3-dimensional subspace. Trajectories of the system are shown coming out from every node. The bold line denotes the 1 -dimensional invariant grid, starting with which the 2 -dimensional grid was constructed.
the absolute value of the symmetrically linearized Jacobian matrix of the system are shown. One can see that the two smallest eigenvalues almost interchange on one of the grid ends. This means that the 1-dimensional 'slow' manifold faces definite problems in this region, it is just not well defined there. In practice, it means that one has to use at least a 2-dimensional grids there.

Figure 7a gives a view of the 2-dimensional invariant grid, constructed for the system, using the 'invariant flag' strategy. The grid was raised starting from the 1-dimensional grid constructed at the previous step. At the first iteration for every node of the initial grid, two nodes (and two edges) were added. The direction of the step was chosen as the direction of the eigenvector of the matrix $A^{s y m}$ (at the point of the node), corresponding to the second 'slowest' direction. The value of the step was chosen to be $\varepsilon=0.05$ in terms of entropic distance. After several Newton's iterations done until convergence was reached, new nodes were added in the direction 'orthogonal' to the 1-dimensional grid. This time it was done by linear extrapolation of the grid on the same step $\varepsilon=0.05$. Once some new nodes become one or several negative coordinates (the grid reaches the boundaries) they were cut off. If a new node has only one edge, connecting it to the grid, it was excluded (since it was impossible to calculate 2-dimensional tangent space for this node). The process continued until the expansion was possible (the ultimate state is when every new node had to be cut off).

The method for calculating tangent vectors for this regular rectangular 2-dimensional grid was chosen to be quite simple. The grid consists of rows, which are co-oriented by construction to the initial 1-dimensional grid, and columns that consist of the adjacent nodes in the neighbouring rows. The direction of the columns corresponds to the second slowest direction along the grid. Then, every row and col-
umn are considered as a 1-dimensional grid, and the corresponding tangent vectors are calculated as it was described before:

$$
g_{\text {row }}\left(x_{k, i}\right)=\left(x_{k, i+1}-x_{k, i-1}\right) /\left\|x_{k, i+1}-x_{k, i-1}\right\|
$$

for the internal nodes and

$$
\begin{aligned}
g_{\text {row }}\left(x_{k, 1}\right) & =\left(x_{k, 1}-x_{k, 2}\right) /\left\|x_{k, 1}-x_{k, 2}\right\|, g_{\text {row }}\left(x_{k, n_{k}}\right) \\
& =\left(x_{k, n_{k}}-x_{k, n_{k}-1}\right) /\left\|x_{k, n_{k}}-x_{k, n_{k}-1}\right\|
\end{aligned}
$$

for the nodes which are close to the grid's edges. Here $x_{k, i}$ denotes the vector of the node in the $k$-th row, $i$-th column; $n_{k}$ is the number of nodes in the $k$-th row. Second tangent vector $g_{\text {col }}\left(x_{k, i}\right)$ is calculated analogously. In practice, it proves convenient to orthogonalize $g_{\text {row }}\left(x_{k, i}\right)$ and $g_{\text {col }}\left(x_{k, i}\right)$.

## 9 Invariant Grid as a Tool for Visualization of Dynamic System Properties

The usual way of dealing with system (1) is to define some initial conditions and solve the equation for a given time interval. This gives us one particular trajectory of the system. Can we have a look at the global picture of all possible trajectories or in other words can we visualize the vector field in $R^{N}$, defined by $\boldsymbol{J}(x)$ ? It would be possible if one had two or three species in system (1). Invariant manifolds and their grid representation allow to do it for higher dimensions, thus they can serve as a data visualization tool. The situation is somewhat close in spirit to data visualization using principal manifolds [for example, see 23] where one uses 2-dimensional manifolds to visualize a finite set of points. Invariant manifolds make it possible to visualize the global system dynamics on the non-linear manifold of slow motions (i.e., in the space which corresponds to the effects observed in a real-life experiment).

In this section we demonstrate global system dynamics visualization on the model hydrogen burning reaction. Since the phase space is 4 -dimensional, it is impossible to visualize the grid in one of the coordinate 3-dimensional views, as it was
done in the previous subsection. To facilitate visualization one can utilize traditional methods of multidimensional data visualization. Here we make use of the principal components analysis [see, for example, 22], which constructs a 3-dimensional linear subspace with maximal dispersion of the orthogonally projected data (grid nodes in our case). In other words, the method of principal components constructs in a multidimensional space a 3-dimensional box such that the grid can be placed maximally tightly inside the box (in the mean square distance meaning). After projection of the grid nodes into this space, we get more or less adequate representation of the 2-dimensional grid embedded into the 6 -dimensional concentrations space (fig. 7b). The disadvantage of the ap-
proach is that the axes now do not bear any explicit physical meaning, they are just some linear combinations of the concentrations.

One attractive feature of 2-dimensional grids is the possibility to use them as a screen, on which one can display different functions $f(\boldsymbol{c})$ defined in the concentrations space. This technology was exploited widely in the non-linear data analysis by the elastic maps method [23,24]. The idea is to 'unfold' the grid on a plane (to present it in the 2-dimensional space, where the nodes form a regular lattice). In other words, we are going to work in the internal coordinates of the grid. In our case, the first internal coordinate (let us call it $s_{1}$ ) corresponds to the direction, co-oriented with the 1-dimensional invariant grid, the sec-


Fig. 8. Two-dimensional invariant grid as a screen for visualizing different functions defined in the concentration space. The coordinate axes are entropic distances (see text for explanations) along the first and the second slowest directions on the grid. The corresponding 1-dimensional invariant grid is denoted by the bold line, the equilibrium is denoted by the square.
ond one (let us call it $s_{2}$ ) corresponds to the second slow direction. By the construction, the coordinate line $s_{2}=0$ corresponds to the 1-dimensional invariant grid. Units of $s_{1}$ and $s_{2}$ is the entropic distance.

Every grid node has two internal coordinates ( $s_{1}, s_{2}$ ) and, simultaneously, corresponds to a vector in the concentration space. This allows us to map any function $f(c)$ from the multidimensional concentration space to the 2-dimensional space of the grid. This mapping is defined in a finite number of points (grid nodes), and can be interpolated (linearly, in the simplest case) between them. Using colouring and isolines one can visualize the values of the function in the neighbourhood of the invariant manifold. This is meaningful, since, by the definition, the system spends most of the time in the vicinity of the invariant manifold; thus, one can visualize the behaviour of the system. As a result of applying this technology, one obtains a set of colour illustrations (a stack of information layers), put onto the grid as a map. This enables applying the whole family of the well-developed methods of working with the stack of information layers, such as the geographical information systems methods.

Briefly, this technique of the visualization is a useful tool for the understanding of dynamical systems. It makes it possible to see many different scenarios of the system behaviour simultaneously, together with different system's characteristics.

Let us use the invariant grids for the model hydrogen burning system as a screen for visualization. The simplest functions to visualize are the coordinates: $c_{i}(\boldsymbol{c})=c_{i}$. In figure 8 we displayed four colourings, corresponding to the four arbitrarily chosen concentration functions (of $\mathrm{H}_{2}, \mathrm{O}, \mathrm{H}$ and OH ; fig. $8 \mathrm{a}-\mathrm{d}$ ). The qualitative conclusion that can be made from the graphs is that, for example, the concentration of $\mathrm{H}_{2}$ practically does not change during the first fast motion (towards the 1-dimensional grid) and, then, gradually changes to the equilibrium value (the $\mathrm{H}_{2}$
coordinate is 'slow'). The $O$ coordinate represents the opposite case; it is the 'fast' coordinate which changes quickly (on the first stage of the motion) to the almost equilibrium value, and it almost does not change after that. Basically, the slopes of the coordinate isolines give some impression of how 'slow' a given concentration is. Figure 8c shows an interesting behaviour of the OH concentration. Close to the 1-dimensional grid it behaves like a'slow coordinate', but there is a region on the map where it has a clear'fast' behaviour (middle bottom of the graph).

The next two functions which can be visualized are the entropy $S$ and the entropy
production $\sigma(\boldsymbol{c})=-d G / d t(\boldsymbol{c})=$ $\sum_{i} \ln \left(c_{i} / c_{i}^{\mathrm{eq}}\right) \dot{c}_{i}$. They are shown in figure 9 a and b .

Finally, we visualize the relation between the relaxation times of the fast motion towards the 2-dimensional grid and the slow motion along it. This is given in figure 9 c . This picture allows us to draw the conclusion that a 2 -dimensional consideration can be appropriate for the system (especially in the 'high $H_{2}$, high $O$ ' region), since the relaxation times 'towards' and 'along' the grid are well separated. One can compare this to figure 9d, where the relation between relaxation times towards and along the 1-dimensional grid is shown.


Fig. 9. Two-dimensional invariant grid as a screen for visualizing different functions defined in the concentration space. The coordinate axes are entropic distances (see text for explanations) along the first and the second slowest directions on the grid. The corresponding 1-dimensional invariant grid is denoted by the bold line, the equilibrium is denoted by the square.

## 10 Invariant Manifolds for Open Systems

10.1 Zero-Order Approximation

Let the initial dissipative system (1) be 'spoiled' by an additional term ('external vector field' $\left.J_{e x}(x, t)\right)$ :

$$
\begin{equation*}
\frac{d x}{d t}=J(x)+J_{e x}(x, t), x \subset U . \tag{53}
\end{equation*}
$$

For this new system the entropy does not increase everywhere. In the new system (53) different dynamic effects are possible, such as a non-uniqueness of stationary states, auto-oscillations, etc. The 'inertial manifold' effect is well-known: solutions of (53) approach some relatively low-dimensional manifold on which all the non-trivial dynamics takes place [2527].

It is natural to expect that the inertial manifold of the system (53) is located somewhere close to the slow manifold of the initial dissipative system (1). This hypothesis has the following basis. Suppose that the vector field $J_{\text {ex }}(x, t)$ is sufficiently small. Let us introduce, for example, a small parameter $\varepsilon>0$, and consider $\varepsilon J_{e x}(x, t)$ instead of $J_{e x}(x, t)$. Let us assume that for system (1) a separation of motions into 'slow' and 'fast' takes place. In this case, there exists such an interval of positive $\varepsilon$ that $\varepsilon J_{e x}(x, t)$ is comparable to $J$ only in a small neighbourhood of the given slow motion manifold of system (1). Outside this neighbourhood, $\varepsilon J_{e x}(x, t)$ is negligibly small in comparison with $J$ and only negligibly influences the motion (for this statement to be true, it is important that system (1) is dissipative and every solution comes in finite time to a small neighbourhood of the given slow manifold).

Precisely this perspective on system (53) allows to exploit slow invariant manifolds constructed for the dissipative system (1) as the ansatz and the zero-order approximation in a construction of the inertial manifold of the open system (53). In the zero-order approximation, the right part of the equation (53) is simply project-
ed onto the tangent space of the slow manifold.

The choice of the projector is determined by the motion separation which was described above: fast motion is taken from the dissipative system (1). A projector which is suitable for all dissipative systems with given entropy function is unique. It is constructed in the following way. Let a point $x \in U$ be defined and some vector space $T$, on which one needs to construct a projection ( $T$ is the tangent space to the slow manifold at the point $x$ ). We introduce the entropic scalar product $\langle\mid\rangle_{x}$ :

$$
\begin{equation*}
\langle a \mid b\rangle_{x}=-\left(a, D_{x}^{2} S(b)\right) \tag{54}
\end{equation*}
$$

Let us consider $T_{0}$ that is a subspace of $T$ and which is annulled by the differential $S$ at the point $x$.

$$
\begin{equation*}
T_{0}=\left\{a \in T \mid D_{x} S(a)=0\right\} \tag{55}
\end{equation*}
$$

If $T_{0}=T$, then the thermodynamic projector is the orthogonal projector on $T$ with respect to the entropic scalar product $\langle\mid\rangle_{x}$. Suppose that $T_{0} \neq T$. Let $e_{g} \in T, e_{g} \perp T_{0}$ with respect to the entropic scalar product $\langle\mid\rangle_{x}$, and $D_{x} S\left(e_{g}\right)=1$. These conditions define vector $e_{g}$ uniquely.

The projector onto $T$ is defined by the formula

$$
\begin{equation*}
P(J)=P_{0}(J)+e_{g} D_{x} s(J) \tag{56}
\end{equation*}
$$

where $P_{0}$ is the orthogonal projector onto $T_{0}$ with respect to the entropic scalar product $\langle\mid\rangle_{x}$. For example, if $T$ a finite-dimensional space, then projector (56) is constructed in the following way. Let $e_{1}, \ldots, e_{n}$ be a basis in $T$, and for definiteness, $D_{x} S\left(e_{1}\right) \neq 0$.

1) Let us construct a system of vectors

$$
\begin{equation*}
b_{i}=E_{i+1}-\lambda_{i} e_{1},(i=1, \ldots, n-1), \tag{57}
\end{equation*}
$$

where $\lambda_{i}=D_{x} S\left(e_{i+1}\right) / D_{x} S\left(e_{1}\right)$, and hence $D_{X} S\left(b_{i}\right)=0$. Thus, $\left\{b_{i}\right\}_{1}^{n-1}$ is a basis in $T_{0}$.
2) Let us orthogonalize $\left\{b_{i}\right\}_{1}^{n-1}$ with respect to the entropic scalar product $\langle\mid\rangle_{x}$ (1). We thus derived an orthonormal with respect to $\langle\mid\rangle_{x}$ basis $\left\{g_{i}\right\}_{1}^{n-1}$ in $T_{0}$.
3) We find $e_{g} \in T$ from the conditions:

$$
\begin{equation*}
\left\langle e_{g} \mid g_{i}\right\rangle_{x}=0,(i=1, \ldots, n-1), D_{x} S\left(e_{g}\right)=1 \tag{58}
\end{equation*}
$$

and, finally we get

$$
\begin{equation*}
P(J)=\sum_{i=1}^{n-1} g_{i}\left\langle g_{i} \mid J\right\rangle_{x}+e_{g} D_{x} S(J) \tag{59}
\end{equation*}
$$

If $D_{\chi} S(T)=0$, then the projector $P$ is simply the orthogonal projector with respect to the $\langle\mid\rangle_{x}$ scalar product. This is possible if $x$ is the global maximum of entropy point (equilibrium).

Then

$$
\begin{equation*}
P(J)=\sum_{i=1}^{n} g_{i}\left\langle g_{i} \mid J\right\rangle_{x},\left\langle g_{i} \mid g_{j}\right\rangle=\delta_{i j} . \tag{60}
\end{equation*}
$$

### 10.2 First-Order Approximation

Thermodynamic projector (56) defines a 'slow and fast motion' duality: if $T$ is the tangent space of the slow motion manifold then $T=\operatorname{im} P$, and $\operatorname{ker} P$ is the plane of fast motions. Let us denote by $P_{x}$ the projector at a point $x$ of a given slow manifold.

The vector field $J_{e x}(x, t)$ can be decomposed in two components:

$$
\begin{equation*}
J_{e x}(x, t)=P_{x} J_{e x}(x, t)+\left(1-P_{x}\right) J_{e x}(x, t) \tag{61}
\end{equation*}
$$

Let us denote $J_{\text {ex } s}=P_{x} J_{\text {ex }}, J_{\text {ex } f}=$ $\left(1-P_{x}\right) J_{e x}$. The slow component $J_{\text {ex } s}$ gives a correction to the motion along the slow manifold. This is a zero-order approximation. The 'fast' component shifts the slow manifold in the fast motions plane. This shift changes $P_{x} J_{e x}$ accordingly. Consideration of this effect gives a first-order approximation. In order to find it, let us rewrite the invariance equation taking $J_{e x}$ into account:

$$
\left\{\begin{array}{l}
\left(1-P_{x}\right)\left(J(x+\delta x)+\varepsilon J_{e x}(x, t)\right)=0  \tag{62}\\
P_{x} \delta x=0
\end{array}\right.
$$

The first iteration of the Newton method subject to incomplete linearization gives:

$$
\begin{align*}
& \left\{\begin{array}{l}
\left(1-P_{x}\right)\left(D_{x} J(\delta x)+\varepsilon J_{e x}(x, t)\right)=0 \\
P_{x} \delta x=0 .
\end{array}\right. \\
& \left(1-P_{x}\right) D_{x} J\left(1-P_{x}\right) J(\delta x)=-\varepsilon J_{e x}(x, t) . \tag{64}
\end{align*}
$$

Thus, we have derived a linear equation in the space ker $P$. The operator $(1-P) D_{x} J$ $(1-P)$ is defined in this space.

Utilization of the self-adjoint linearization instead of the traditional linearization $D_{X} J$ operator considerably simplifies solving and studying equation (64). It is necessary to take into account here that the projector $P$ is a sum of the orthogonal projector with respect to the $\langle\mid\rangle_{x}$ scalar product and a projector of rank one.

Assume that the first-order approximation equation (64) has been solved and the following function has been found:

$$
\begin{equation*}
\delta_{1} x\left(x, \varepsilon J_{e x f}\right)=-\left[\left(1-P_{x}\right) D_{x} J\left(1-p_{x}\right)\right]^{-1} \varepsilon \delta_{e x x}, \tag{65}
\end{equation*}
$$

where $D_{x} J$ is either the differential of $J$ or symmetrized differential of $J$ (20).

Let $x$ be a point on the initial slow manifold. At the point $x+\delta x\left(x, \varepsilon J_{e x} f\right)$ the right-hand side of equation (53) in the first-order approximation is given by

$$
\begin{equation*}
J(x)+\varepsilon J_{e x}(x, t)+D_{x} J\left(\delta x\left(x, \varepsilon J_{e x} f\right)\right) . \tag{66}
\end{equation*}
$$

Due to the first-order approximation (66), the motion of a point projection onto the manifold is given by the following equation

$$
\begin{align*}
\frac{d x}{d t}= & P_{x}\left(J(x)+\varepsilon J_{e x}(x, t)+\right. \\
& \left.D_{x} J\left(\delta x\left(x, \varepsilon J_{e x f}(x, t)\right)\right)\right) . \tag{67}
\end{align*}
$$

Note that, in equation (67), the vector field $J(x)$ enters only in the form of a projection, $P_{x} J(x)$. For the invariant slow manifold it holds $P_{x} J(x)=J(x)$, but actually we always deal with approximately invariant manifolds; hence, it is necessary to use the projection $P_{x} J$ instead of $J$ in (67).

Remark. The notion 'projection of a point onto the manifold' needs to be specified. For every point $x$ of the slow invariant manifold $M$ both the thermodynamic pro-
jector $P_{x}(56)$ and the fast motions plane $\operatorname{ker} P_{x}$ are defined. Let us define a projector $\Pi$ of some neighbourhood of $M$ onto $M$ in the following way:

$$
\begin{equation*}
\Pi(z)=x \text {, if } P_{x}(z-x)=0 . \tag{68}
\end{equation*}
$$

Qualitatively, it means that $z$, after all fast motions took place, comes into a small neighbourhood of $x$. Operation (56) is defined uniquely in some small neighbourhood of the manifold $M$.

A derivation of slow motion equations requires not only an assumption that $\varepsilon J_{e x}$ is small but it must be slow as well: $\frac{d}{d t}\left(\varepsilon J_{e x}\right)$ must be small too.

Further approximations for slow motions of system (53) can be obtained, taking into account the time derivatives of $J_{e x}$. This is an alternative to the usage of the projection operator methods [28].

## 11 Conclusion

In this paper we presented a method for reducing the complexity in complex chemical reaction networks using a consistent approach of constructing invariant manifold for the system of kinetic equations. The method is applicable to the class of dissipative systems (with Lyapounov function) and can be extended to the case of open systems as well.

An attractive feature of the approach is its clear geometrical interpretation. The geometrical approach becomes more and more popular in applied model reduction: one constructs a slow approximate invariant manifold, and dynamical equations on this manifold instead of an approximation of solutions to the initial equations. After that, the equations on the slow manifold can be studied separately, as well as the fast motion to this manifold (the initial layer problem [29]).

The notion of invariant grid may be useful beyond the chemical kinetics. This discrete invariant object can serve as a representation of approximate slow invariant manifold, and as a screen (a map) for visualization of different functions and prop-
erties. The problem of the grid correction is fully decomposed into the problems of the grid's node correction which makes it open to effective parallel implementations.

The next step should be the implementation of the MIG for investigation of highdimensional systems 'kinetics + transport'. The asymptotic analysis of the methods of analytic continuation the manifold from the grid should lead to further development of these methods and modifications of the Carleman formula.

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