# $\mathcal{C}^{r}$-Lohner algorithm 

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#### Abstract

We present a Lohner type algorithm for the computation of rigorous bounds for the solutions of ordinary differential equations and its derivatives with respect to the initial conditions up to an arbitrary order.


Keywords: rigorous integration of ODEs, variational equations

## 1 Introduction

This paper is a sequel to [Z]. We present here a Lohner-type algorithm for the computation of rigorous enclosures of the partial derivatives with respect to the initial conditions up to an arbitrary order $r$ of the flow induced by an autonomous ODE, hence the name the $\mathcal{C}^{r}$-Lohner algorithm. Let $r$ be a positive integer, then by the $\mathcal{C}^{r}$-algorithm we will mean a routine which gives rigorous estimates for the partial derivatives with respect to the initial conditions up to the order $r$ and by the $\mathcal{C}^{r}$-computations we mean an application of the $\mathcal{C}^{r}$ algorithm.

Our main motivation for the development of the $\mathcal{C}^{r}$-algorithm was a desire to provide a tool, which will considerably extend the possibilities of computer assisted proofs in the dynamics of ODEs, which require rigorous bounds on orbits. Till now most of such proofs have used the topological conditions (see for example [HZHT, MM, GZ, Z1]) and additionally the conditions on the first derivatives with respect to the initial conditions (see for example [RNS, T, WZ, $\mathrm{KZ}]$ ), hence it required the $\mathcal{C}^{0}$ - and $\mathcal{C}^{1}$-computations, respectively. The spectrum of problems addressed includes the questions of the existence of periodic orbits and their local uniqueness, the existence of symbolic dynamics, the existence of hyperbolic invariants sets, the existence of homo- and heteroclinic orbits. To address other phenomena, like the bifurcations of periodic orbits, the route to chaos, invariant tori through the KAM theory one needs the knowledge of partial derivatives with respect to the initial conditions of the higher order.

In principle, one can think that a good rigorous ODE solver should be enough. Namely, to compute the partial derivatives of the flow induced by

$$
\begin{equation*}
x^{\prime}=f(x), \quad x \in \mathbb{R}^{n} \tag{1}
\end{equation*}
$$

[^0]it is enough to rigorously integrate the system of variational equations obtained by the formal differentiation of (1) with respect to the initial conditions. For example for $r=2$ we have the following system
\[

$$
\begin{align*}
x^{\prime} & =f(x)  \tag{2}\\
\frac{d}{d t} V_{i j}(t) & =\sum_{s=1}^{n} \frac{\partial f_{i}}{\partial x_{s}}(x) V_{s j}(t)  \tag{3}\\
\frac{d}{d t} H_{i j k}(t) & =\sum_{s, r=1}^{n} \frac{\partial^{2} f_{i}}{\partial x_{s} \partial x_{r}}(x) V_{r k}(t) V_{s j}(t)+\sum_{s=1}^{n} \frac{\partial f_{i}}{\partial x_{s}}(x) H_{s j k}(x) \tag{4}
\end{align*}
$$
\]

with the initial conditions

$$
\begin{equation*}
x(0)=x_{0}, \quad V(0)=\mathrm{Id}, \quad H_{i j k}(0)=0, \quad i, j, k=1, \ldots, n . \tag{5}
\end{equation*}
$$

It is well known that if by $\varphi\left(t, x_{0}\right)$ we denote the (local) flow induced by (1), then

$$
\begin{aligned}
\frac{\partial \varphi_{i}}{\partial x_{j}}\left(t, x_{0}\right) & =V_{i, j}(t), \\
\frac{\partial^{2} \varphi_{i}}{\partial x_{j} \partial x_{k}}\left(t, x_{0}\right) & =H_{i j k}(t) .
\end{aligned}
$$

Analogous statements are true for the higher order partial derivatives with respect to the initial conditions.

Remark 1 The variational equations up to an arbitrary order might be generated automatically by means of the automatic differentiation [G, Ra]. The main reason for which we discuss in this paper an explicit compact formula for the equations of variations (see (3,4) and Section 2) is to explain a method for the generation of the rough enclosure for the solution of higher order variational equations. In the practical implementation the use of any compact formulas for variational equations can be avoided.

It turns out that a straightforward application of any rigorous ODE solver to the system of variational equations $(2-4)$ is very inefficient. Namely, it totally ignores the structure of the system and sometimes it leads to a very poor performance and unnecessary long computation times (see Section 3.1 for more discussion and Section 7 for results of our tests).

Our algorithm is a modification of the Lohner algorithm [Lo], which takes into account the structure of variational equations (2-4). Basically it consists of the Taylor method, a heuristic routine for a priori bounds for the solution of (2-4) during a time step and a Lohner-type control of the wrapping effect, which is done separately for $x$ and the partial derivatives with respect to the initial conditions (the variables $V$ and $H$ in (3,4)).

The proposed algorithm has been already successfully applied to several problems. In [KWZ] a computer assisted proof of the existence of the cocooning
cascade of heteroclinic tangencies for the Michelson system [Mi] was given. This proof required the $\mathcal{C}^{2}$ computations. That time we had a special implementation of the $\mathcal{C}^{2}$ algorithm, only.

In [WZ2] the method for proving the existence of quadratic homoclinic tangencies for maps is proposed. An application of the method to a Poincaré map for forced-damped pendulum system required $\mathcal{C}^{2}$ computations. In [WZ3] an application of the $\mathcal{C}^{3}$ algorithm to rigorous verification of period doubling bifurcations for the Rössler system [R] is presented.

In [Wi] $\mathcal{C}^{3}$ and $\mathcal{C}^{5}$ computations were used to prove the existence of invariant tori around some elliptic periodic orbits for hamiltonian and reversible systems. The approach based on the classical KAM theorem for twist maps on the plane.

We believe, the proposed algorithm has a wide spectrum of other applications.

## 2 Faá di Bruno formula.

To effectively deal with the formulas involving the partial derivatives of the composition of maps we will use extensively a notation of multiindices, multipointers and submultipointers throughout the paper. In particular, when used, the variational equations can be written in a compact form.

### 2.1 Multiindices

By $\mathbb{N}$ we will denote the set of nonnegative integers, i.e. $\mathbb{N}=\{0,1,2, \ldots\}$.
Definition 2 An element $\tau \in \mathbb{N}^{n}$ will be called a multiindex.
For a sequence $\alpha=\left(\alpha_{1}, \ldots, \alpha_{n}\right) \in \mathbb{N}^{n}$ and a vector $x=\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{R}^{n}$ we set

1. $|\alpha|=\alpha_{1}+\cdots+\alpha_{n}$,
2. $\alpha!=\alpha_{1} \cdot \alpha_{2} \cdots \alpha_{n}$,
3. $x^{\alpha}=\left(x_{1}^{\alpha_{1}}, \ldots, x_{n}^{\alpha_{n}}\right)$.

By $e_{i}^{n} \in \mathbb{N}^{n}$ we will denote

$$
e_{i}^{n}=(0,0, \ldots, 0, \overbrace{1}^{i}, 0, \ldots, 0,0) .
$$

We will drop the index $n$ (the dimension) in the symbol $e_{i}^{n}$ when it is obvious from the context.

Put $\mathbb{N}_{p}^{n}:=\left\{a \in \mathbb{N}^{n}:|a|=p\right\}$. For $\delta=\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathbb{N}^{n_{1}} \times \cdots \times \mathbb{N}^{n_{k}}$ we set

$$
|\delta|=\sum_{i=1}^{k}\left|\delta_{i}\right| .
$$

Let $f=\left(f_{1}, \ldots, f_{m}\right): \mathbb{R}^{n} \rightarrow \mathbb{R}^{m}$ be sufficiently smooth. For $\alpha \in \mathbb{N}^{n}$ we set

1. $D^{\alpha} f_{i}=\frac{\partial^{|\alpha|} f_{i}}{\partial x_{1}^{\alpha_{1}} \cdots \partial x_{n}^{\alpha_{n}}}$,
2. $D^{\alpha} f=\left(D^{\alpha} f_{1}, D^{\alpha} f_{2}, \ldots, D^{\alpha} f_{m}\right)$.

For a function $f: \mathbb{R} \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ by $D^{\alpha} f_{i}(t, x)$ we will denote $D^{\alpha} f_{i}(t, \cdot)(x)$ and similarly

$$
D^{\alpha} f(t, x)=\left(D^{\alpha} f_{1}(t, x), \ldots, D^{\alpha} f_{n}(t, x)\right)
$$

This convention means that $D^{\alpha}$ always acts on $x$-variables.

### 2.2 Multipointers

For a fixed $n>0$ and $p>0$ we define

$$
\begin{aligned}
\mathcal{N}_{p}^{n} & =\left\{\left(a_{1}, a_{2}, \ldots, a_{p}\right) \in \mathbb{N}^{p}: 1 \leq a_{1} \leq \cdots \leq a_{p} \leq n\right\} \\
\mathcal{N}=\mathcal{N}^{n} & =\bigcup_{p=1}^{\infty} \mathcal{N}_{p}^{n}
\end{aligned}
$$

Definition 3 An element of $\mathcal{N}^{n}$ will be called a multipointer.
Remark 4 A function

$$
\begin{equation*}
\Lambda: \mathcal{N}_{p}^{n} \ni\left(a_{1}, \ldots, a_{p}\right) \rightarrow \sum_{i=1}^{p} e_{a_{i}}^{n} \in \mathbb{N}_{p}^{n} \tag{6}
\end{equation*}
$$

is a bijection.
Let $f=\left(f_{1}, \ldots, f_{m}\right): \mathbb{R}^{n} \rightarrow \mathbb{R}^{m}$ be a sufficiently smooth. For $a \in \mathcal{N}_{p}^{n}$ we set

1. $D_{a} f_{i}=\frac{\partial^{p} f_{i}}{\partial x_{a_{1}} \ldots \partial x_{a_{p}}}$,
2. $D_{a} f=\left(D_{a} f_{1}, \ldots, D_{a} f_{m}\right)$.

For a function $f: \mathbb{R} \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ by $D_{a} f_{i}(t, x)$ we will denote $D_{a} f_{i}(t, \cdot)(x)$. In the light of the above notations $D_{\alpha} f=D^{\Lambda(\alpha)} f$.

For $a=\left(a_{1}, a_{2}, \ldots, a_{n}\right) \in \mathbb{N}_{p}^{n}$ and $b=\left(b_{1}, b_{2}, \ldots, b_{n}\right) \in \mathbb{N}_{q}^{n}$ we define

$$
a+b=\left(a_{1}+b_{1}, \ldots, a_{n}+b_{n}\right) \in \mathbb{N}_{p+q}^{n} .
$$

For $\alpha \in \mathcal{N}_{p}^{n}$ and $\beta \in \mathcal{N}_{q}^{n}$ we define

$$
\alpha+\beta=\Lambda^{-1}(\Lambda(\alpha)+\Lambda(\beta)) \in \mathcal{N}_{p+q}^{n}
$$

By $\leq$ we will denote a linear order (the lexicographical order) in $\mathcal{N}$ defined in the following way. For $a \in \mathcal{N}_{p}^{n}$ and $b \in \mathcal{N}_{q}^{n}$

$$
(a \leq b) \Longleftrightarrow\left\{\begin{array}{l}
\text { either } \exists i, i \leq p, i \leq q, a_{i}<b_{i} \text { and } a_{j}=b_{j} \text { for } j<i  \tag{7}\\
\text { or } p \leq q \text { and } a_{i}=b_{i} \text { for } i=1, \ldots, p
\end{array}\right.
$$

Definition 5 For $k \leq p$ we set

$$
\begin{equation*}
\mathcal{N}^{p}(k):=\left\{\left(\delta_{1}, \ldots, \delta_{k}\right) \in\left(\mathcal{N}^{p}\right)^{k}: \delta_{1} \leq \cdots \leq \delta_{k}, \delta_{1}+\cdots+\delta_{k}=(1,2, \ldots, p)\right\} \tag{8}
\end{equation*}
$$

We will use $\mathcal{N}^{p}(k)$ extensively in the next section. Its will be used to label terms in $D^{\alpha} f_{i}(\varphi(t, x))$. Observe that for $p>0$

$$
\begin{array}{r}
\mathcal{N}^{p}(1)=\{(1,2, \ldots, p)\} \\
\mathcal{N}^{p}(p)=\{((1),(2), \ldots,(p))\} .
\end{array}
$$

One can construct all elements of $\mathcal{N}^{p}(k)$ using the following recursive procedure. From the definition of $\mathcal{N}^{p}(k)$ it follows that if $\left(\delta_{1}, \ldots, \delta_{m-1}\right) \in \mathcal{N}^{p-1}(m-1)$ then $\left(\delta_{1}, \ldots, \delta_{m-1},(p)\right) \in \mathcal{N}^{p}(m)$ (notice that order is preserved). Similarly, if $\left(\delta_{1}, \ldots, \delta_{m}\right) \in \mathcal{N}^{p-1}(m)$ then

$$
\left(\delta_{1}, \ldots, \delta_{s-1}, \delta_{s}+(p), \delta_{s+1}, \ldots, \delta_{m}\right) \in \mathcal{N}^{p}(m)
$$

and again order of elements is preserved. Hence, for $p>2$ and $1<k<p$ we have $\mathcal{N}^{p}(k)=A \cup B$ where

$$
\begin{align*}
A & =\left\{\left(\delta_{1}, \ldots, \delta_{k-1},(p)\right):\left(\delta_{1}, \ldots, \delta_{k-1}\right) \in \mathcal{N}^{p-1}(k-1)\right\} \\
B & =\bigcup_{s=1}^{k}\left\{\left(\delta_{1}, \ldots, \delta_{s-1}, \delta_{s}+(p), \delta_{s+1}, \ldots, \delta_{k}\right):\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p-1}(k)\right\} \tag{9}
\end{align*}
$$

and the sets $A$ and $B$ are disjoint.
Another way to generate all elements of $\mathcal{N}^{p}(k)$ can be described as follows

- decompose the set $\{1,2, \ldots, p\}$ into $k$ nonempty and disjoints sets $\Delta_{i}$, $i=1, \ldots, k$,
- we sort each $\Delta_{i}$ and permute $\Delta_{i}$ 's to obtain $\min \left(\Delta_{1}\right)<\min \left(\Delta_{2}\right)<\cdots<$ $\min \left(\Delta_{k}\right)$,
- we define $\delta_{i}$ to be an ordered set consisting of all elements of $\Delta_{i}$ for $i=1, \ldots, k$.

Definition 6 For an arbitrary $a \in \mathcal{N}_{p}^{n}$ and $\delta \in \mathcal{N}_{k}^{p}$ such that $k \leq p$ we define a submultipointer $a_{\delta} \in \mathcal{N}_{k}^{n}$ by $\left(a_{\delta}\right)_{i}=a_{\delta_{i}}$ for $i=1, \ldots, k$, which can be expressed using $\Lambda$ as follows

$$
a_{\delta}:=\Lambda^{-1}\left(\sum_{i=1}^{k} e_{a_{\delta_{i}}}^{n}\right) \in \mathcal{N}_{k}^{n}
$$

### 2.3 The variational equations

Consider an ODE $x^{\prime}=f(x)$ where $f$ is $\mathcal{C}^{K+1}$. Let $\varphi: D \subset \mathbb{R} \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ be a local dynamical system induced by $x^{\prime}=f(x)$. It is well known, that $\varphi \in \mathcal{C}^{K}$ and one can derive the equations for partial derivatives of $\varphi$ by differentiating
equation $\frac{\partial \varphi}{\partial t}(t, x)=f(\varphi(t, x))$ with respect to the initial condition $x$. As a result we obtain a system of so-called equations for variations, the size of which depends on the order $r$ of the partial derivatives we intend to compute. An example of such system for $r=2$ is given by (2-4) with the initial conditions given by (5).

The equations for the higher order partial derivatives written in a compact form using multipointers and multiindices are given by the Faá di Bruno formula.

Lemma 7 ([H], Faá di Bruno formula) For any p-times continuously differentiable functions $f, g: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ and $a \in \mathcal{N}_{p}^{n}$ we have

$$
\begin{equation*}
D_{a}(f \circ g)=\sum_{k=1}^{p} \sum_{i_{1}, \ldots, i_{k}=1}^{n}\left(D^{e_{i_{1}}+\cdots+e_{i_{k}}} f_{i}\right) \circ g \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p}(k)} \prod_{j=1}^{k} D_{a_{\delta_{j}}} g_{i_{j}} . \tag{10}
\end{equation*}
$$

Proof: In the proof the functions $D^{e_{i_{1}}+\cdots+e_{i_{k}}} f_{i}$ are always evaluated at $g(x)$, and various partial derivatives of $g$ are always evaluated at $x$, therefore the arguments will be always dropped to simplify formulae.

Put $F=f \circ g$. We prove the lemma by induction on $p=|a|$. If $p=1$ then $a=(c)$ for some $c \in\{1, \ldots, n\}$ and (15) becomes

$$
D_{(c)} F=\sum_{s=1}^{n} \frac{\partial f_{i}}{\partial x_{s}} \frac{\partial g_{s}}{\partial x_{c}}=\sum_{s=1}^{n} D^{e_{s}} f_{i} \cdot D_{(c)} g_{s}
$$

Assume (15) holds true for $p-1, p>1$. Let us fix $a \in \mathcal{N}_{p}^{n}$. We have $a=b+(c)$, where $b=\left(a_{1}, \ldots, a_{p-1}\right) \in \mathcal{N}_{p-1}^{n}$ and $c=a_{p}$. Since (15) is satisfied for $p-1$, therefore we have

$$
\begin{aligned}
D_{a} F_{i} & =D_{(c)}\left(D_{b} F_{i}\right) \\
& =D_{(c)}\left(\sum_{k=1}^{p-1} \sum_{\substack{i_{1}, \ldots, i_{k}=1 \\
\beta:=e_{i_{1}}+\cdots+e_{i_{k}}}}^{n} D^{\beta} f_{i} \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p-1}(k)} \prod_{j=1}^{k} D_{b_{\delta_{j}}} g_{i_{j}}\right) \\
& =\sum_{k=1}^{p-1} \sum_{\substack{i_{1}, \ldots, i_{k+1}=1 \\
\beta:=e_{i_{1}}+\cdots+e_{i_{k+1}}}}^{n} D^{\beta} f_{i} \cdot D_{(c)} g_{i_{k+1}} \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p-1}(k)} \prod_{j=1}^{k} D_{b_{\delta_{j}}} g_{i_{j}} \\
& +\sum_{k=1}^{p-1} \sum_{\substack{i_{1}, \ldots, i_{k}=1 \\
\beta:=e_{i_{1}}+\cdots+e_{i_{k}}}}^{n} D^{\beta} f_{i} \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p-1}(k)} \sum_{s=1}^{k} D_{b_{\delta_{s}}+(c)} g_{i_{s}} \prod_{\substack{j=1, j \neq s}}^{k} D_{b_{\delta_{j}}} g_{i_{j}} .
\end{aligned}
$$

For $k=1, \ldots, p$ we set

$$
\begin{equation*}
T_{k}:=\sum_{i_{1}, \ldots, i_{k}=1}^{n} D^{e_{i_{1}}+\cdots+e_{i_{k}}} f_{i} \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p}(k)} \prod_{j=1}^{k} D_{a_{\delta_{j}}} g_{i_{j}} . \tag{11}
\end{equation*}
$$

Now our goal is to prove that:

$$
\begin{equation*}
D_{a} F_{i}=\sum_{k=1}^{p} T_{k} \tag{12}
\end{equation*}
$$

Our strategy of proof is as follows. We will define $S_{1}, \ldots, S_{p}$, such that

$$
\begin{equation*}
D_{a} F_{i}=\sum_{k=1}^{p} S_{k} \tag{13}
\end{equation*}
$$

and we will show that $S_{i}=T_{i}$ for $i=1, \ldots, p$.
We set

$$
\begin{aligned}
& S_{1}=\sum_{k=1}^{1} \sum_{\substack{i_{1}, \ldots, i_{k}=1 \\
\beta:=e_{i_{1}}+\cdots+e_{i_{k}}}}^{n} D^{\beta} f_{i} \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p-1}(k)} \sum_{s=1}^{k} D_{b_{\delta_{s}}+(c)} g_{i_{s}} \prod_{\substack{j=1, j \neq s}}^{k} D_{b_{\delta_{j}}} g_{i_{j}}, \\
& S_{p}=\sum_{k=p-1}^{p-1} \sum_{\substack{i_{1}, \ldots, i_{k+1}=1 \\
\beta:=e_{i_{1}}+\cdots+e_{i_{k+1}}}}^{n} D^{\beta} f_{i} \cdot D_{(c)} g_{i_{k+1}} \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p-1}(k)} \prod_{j=1}^{k} D_{b_{\delta_{j}}} g_{i_{j}} .
\end{aligned}
$$

For $m=2,3, \ldots, p-1$ we set

$$
\begin{aligned}
S_{m} & =\sum_{k=m-1}^{m-1} \sum_{\substack{i_{1}, \ldots, i_{k+1}=1 \\
\beta:=e_{i_{1}}+\cdots+e_{i_{k+1}}}}^{n} D^{\beta} f_{i} \cdot D_{(c)} g_{i_{k+1}} \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p-1}(k)} \prod_{j=1}^{k} D_{b_{\delta_{j}}} g_{i_{j}} \\
& +\sum_{k=m}^{m} \sum_{\substack{i_{1}, \ldots, i_{k}=1 \\
\beta:=e_{i_{1}}+\cdots+e_{i_{k}}}}^{n} D^{\beta} f_{i} \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p-1}(k)} \sum_{s=1}^{k} D_{b_{\delta_{s}}+(c)} g_{i_{s}} \prod_{\substack{j=1, j \neq s}}^{k} D_{b_{\delta_{j}}} g_{i_{j}} .
\end{aligned}
$$

There remains for us to show that $S_{i}=T_{i}$ for $i=1, \ldots, p$. Consider first $i=1$. Recall that $\mathcal{N}^{p-1}(1)=\{(1,2, \ldots, p-1)\}$, hence

$$
S_{1}=\sum_{s=1}^{n} D^{e_{s}} f_{i} \cdot D_{b+(c)} g_{s}=\sum_{s=1}^{n} D^{e_{s}} f_{i} \cdot D_{a} g_{s}
$$

Therefore

$$
\begin{equation*}
S_{1}=T_{1} . \tag{14}
\end{equation*}
$$

Consider now $i=p$. For an arbitrary $s>0 \mathcal{N}^{s}(s)$ contains only one element $((1),(2), \ldots,(s))$. Therefore we obtain

$$
\begin{aligned}
S_{p} & =\sum_{i_{1}, \ldots, i_{p}=1}^{n} D^{e_{i_{1}}+\cdots+e_{i_{p}}} f_{i} \cdot D_{(c)} g_{i_{p}} \sum_{\left(\delta_{1}, \ldots, \delta_{p-1}\right) \in \mathcal{N}^{p-1}(p-1)} \prod_{j=1}^{p-1} D_{b_{\delta_{j}}} g_{i_{j}} \\
& =\sum_{i_{1}, \ldots, i_{p}=1}^{n} D^{e_{i_{1}}+\cdots+e_{i_{p}}} f_{i} \cdot D_{(c)} g_{i_{p}} \prod_{j=1}^{p-1} D_{b_{j}} g_{i_{j}} .
\end{aligned}
$$

Since $a=b+(c)$, where $c=\left(a_{p}\right)$, we have

$$
\begin{aligned}
S_{p} & =\sum_{i_{1}, \ldots, i_{p}=1}^{n} D^{e_{i_{1}}+\cdots+e_{i_{p}}} f_{i} \prod_{j=1}^{p} D_{a_{j}} g_{i_{j}} \\
& =\sum_{i_{1}, \ldots, i_{p}=1}^{n} D^{e_{i_{1}}+\cdots+e_{i_{p}}} f_{i} \sum_{\left(\delta_{1}, \ldots, \delta_{p}\right) \in \mathcal{N}^{p}(p)} \prod_{j=1}^{p} D_{a_{\delta_{j}}} g_{i_{j}}=T_{p} .
\end{aligned}
$$

Consider now $m=2,3, \ldots, p-1$. We have

$$
\begin{aligned}
S_{m} & =\sum_{i_{1}, \ldots, i_{m}=1}^{n} D^{e_{i_{1}}+\cdots+e_{i_{m}}} f_{i} \cdot D_{(c)} g_{i_{m}} \sum_{\left(\delta_{1}, \ldots, \delta_{m-1}\right) \in \mathcal{N}^{p-1}(m-1)} \prod_{j=1}^{m-1} D_{b_{\delta_{j}}} g_{i_{j}} \\
& +\sum_{i_{1}, \ldots, i_{m}=1}^{n} D^{e_{i_{1}}+\cdots+e_{i_{m}}} f_{i} \sum_{\left(\delta_{1}, \ldots, \delta_{m}\right) \in \mathcal{N}^{p-1}(m)} \sum_{s=1}^{m} D_{b_{\delta_{s}}+(c)} g_{i_{s}} \prod_{\substack{j=1, j \neq s}}^{m} D_{b_{\delta_{j}}} g_{i_{j}}
\end{aligned}
$$

Using decomposition $\mathcal{N}^{p}(m)=A \cup B$ as in (9) we obtain

$$
\begin{aligned}
S_{m} & =\sum_{i_{1}, \ldots, i_{m}=1}^{n} D^{e_{i_{1}}+\cdots+e_{i_{m}}} f_{i} \sum_{\left(\delta_{1}, \ldots, \delta_{m-1}, \delta_{m}=(p)\right) \in A} \prod_{j=1}^{m} D_{a_{\delta_{j}}} g_{i_{j}} \\
& +\sum_{i_{1}, \ldots, i_{m}=1}^{n} D^{e_{i_{1}}+\cdots+e_{i_{m}}} f_{i} \sum_{\left(\delta_{1}, \ldots, \delta_{m}\right) \in B} \prod_{j=1}^{m} D_{a_{\delta_{j}}} g_{i_{j}} \\
& =\sum_{i_{1}, \ldots, i_{m}=1}^{n} D^{e_{i_{1}}+\cdots+e_{i_{m}}} f_{i} \sum_{\left(\delta_{1}, \ldots, \delta_{m}\right) \in \mathcal{N}^{p}(m)} \prod_{j=1}^{m} D_{a_{\delta_{j}}} g_{i_{j}}=T_{m}
\end{aligned}
$$

We have shown that $T_{i}=S_{i}$ for $i=1, \ldots, p$. This finishes the proof.
From the above lemma we have immediately
Lemma 8 Assume $f \in \mathcal{C}^{r+1}$ and let $\varphi: \mathbb{R} \times \mathbb{R}^{n} \rightarrow \rightarrow \mathbb{R}^{n}$ be a local dynamical system induced by $x^{\prime}=f(x)$. Then for $a \in \mathcal{N}_{p}^{n}$ such that $p \leq r$ holds

$$
\begin{equation*}
\frac{d}{d t} D_{a} \varphi_{i}=\sum_{k=1}^{p} \sum_{i_{1}, \ldots, i_{k}=1}^{n}\left(D^{e_{i_{1}}+\cdots+e_{i_{k}}} f_{i}\right) \circ \varphi \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p}(k)} \prod_{j=1}^{k} D_{a_{\delta_{j}}} \varphi_{i_{j}} \tag{15}
\end{equation*}
$$

for $i=1, \ldots, n$.
Formula (15) could be seen as a direct application of the chain rule for composition of multivariate power series. Using the automatic differentiation tools [JZ, G, Ra] one can efficiently nonrigorously integrate ODE's together with higher order variational equations by means of floating point arithmetic.

The main goal of this paper is to present an efficient rigorous solver for higher order variational equations which takes into account the structure of the equations and the wrapping effect.

## $3 \mathcal{C}^{r}$-Lohner algorithm

### 3.1 Why one needs an $\mathcal{C}^{r}$-algorithm?

There exist in the literature several effective algorithms for the computation of the rigorous bounds for the solutions of ordinary differential equations, including the Lohner method [Lo], the Hermite-Obreschkoff algorithm [NJ] or the Taylor model $[\mathrm{BM}]$. For the $\mathcal{C}^{r}$-computations the number of equations to solve is equal to $n\binom{n+r}{n}$ hence, even for $r=1$ the direct application of such algorithms to the equations for variations (16) leads to the integration in the high dimensional space and is usually inefficient. Let us recall after [Z, Sec. 6] the basic reason for this. In order to have a good control over the expansion rate of the set of the initial conditions during a time step these algorithms, while being $\mathcal{C}^{0}$, are $\mathcal{C}^{1}$ 'internally'(or higher for the Taylor models), because they solve non-rigorously the equations for $\left(\frac{\partial \varphi}{\partial x}\right)$ - the variational matrix of the flow. This effectively squares the dimension of phase space of the equation and impacts heavily on the computation time. But as it was observed in [Z] the equations for the partial derivatives of the flow can be seen as the non-autonomous and nonhomogenous linear system of equations, therefore we do not need any additional equations for variations for them. As a result the dimension of the effective phase space for our $\mathcal{C}^{r}$-algorithm is given by $n\binom{n+r}{n}$ instead of the square of this number.

Another important aspect of the proposed algorithm is the fact that the Lohner-type control of the wrapping effect is done separately for $x$-variables and variables $D_{a} \varphi$. This feature is not present in the naive application of $\mathcal{C}^{0}$ algorithm to the system of variational equations and it turns out that this often practically switches off the control of the wrapping effect on $x$-variables, as various choices used in this control become dominated by the $D_{a} \varphi$-variables.

In Section 7 we will give detailed comparison of the $\mathcal{C}^{0}$-solver applied to the equations of variations and our $\mathcal{C}^{r}$-solver.

### 3.2 An outline of the algorithm

Let us fix $r \leq K$ and consider the following system of differential equations

$$
\left\{\begin{align*}
\frac{d}{d t} \varphi & =f \circ \varphi  \tag{16}\\
\frac{d}{d t} D_{a} \varphi & =\sum_{k=1}^{d} \sum_{i_{1}, \ldots, i_{k}=1}^{n}\left(D^{e_{i_{1}}+\cdots+e_{i_{k}}} f\right) \circ \varphi \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{d}(k)} \prod_{j=1}^{k} D_{a_{\delta_{j}}} \varphi_{i_{j}}
\end{align*}\right.
$$

for all $a \in \mathcal{N}_{d}^{n}, d=1, \ldots, r$.
The initial conditions for (16) are

$$
\begin{cases}\varphi\left(0, x_{0}\right) & \in\left[x_{0}\right] \subset \mathbb{R}^{n},  \tag{17}\\ D \varphi\left(0, x_{0}\right) & =\mathrm{Id}, \\ D_{a} \varphi\left(0, x_{0}\right) & =0, \quad \text { for } a \in \mathcal{N}_{2}^{n} \cup \ldots \cup \mathcal{N}_{r}^{n} .\end{cases}
$$

In the sequel we will use the following notations:

- if a solution of system (16) is defined for $t>0$ and some $x_{0} \in \mathbb{R}^{n}$, then for $a \in \mathcal{N}$ by $V_{a}\left(t, x_{0}\right)$ we denote $D_{a} \varphi\left(t, x_{0}\right)$,
- for $\left[x_{0}\right] \subset \mathbb{R}^{n}$ by $\left[V_{a}\left(t,\left[x_{0}\right]\right)\right]$ we will denote a set for which we have $V_{a}\left(t,\left[x_{0}\right]\right) \subset\left[V_{a}\left(t,\left[x_{0}\right]\right)\right]$. This set is obtained using an rigorous numerical routine described below.

The $\mathcal{C}^{r}$-Lohner algorithm is a modification of the $\mathcal{C}^{1}$-Lohner algorithm [Z]. One step of the $\mathcal{C}^{r}$-Lohner is a shift along the trajectory of the system (16) with the following input and output data

## Input data:

- $t_{k}$ - a current time,
- $h_{k}$ - a time step,
- $\left[x_{k}\right] \subset \mathbb{R}^{n}$, such that $\varphi\left(t_{k},\left[x_{0}\right]\right) \subset\left[x_{k}\right]$,
- $\left[V_{k, a}\right]=\left[V_{k, a}\left(t_{k},\left[x_{0}\right]\right)\right] \subset \mathbb{R}^{n}$, such that $D_{a} \varphi\left(t_{k},\left[x_{0}\right]\right) \subset\left[V_{k, a}\right]$ for $a \in$ $\mathcal{N}_{1}^{n} \cup \ldots \cup \mathcal{N}_{r}^{n}$.


## Output data:

- $t_{k+1}=t_{k}+h_{k}$ - a new current time,
- $\left[x_{k+1}\right] \subset \mathbb{R}^{n}$, such that $\varphi\left(t_{k+1},\left[x_{0}\right]\right) \subset\left[x_{k+1}\right]$,
- $\left[V_{k+1, a}\right]=\left[V_{k+1, a}\left(t_{k+1},\left[x_{0}\right]\right)\right] \subset \mathbb{R}^{n}$, such that $D_{a} \varphi\left(t_{k+1},\left[x_{0}\right]\right) \subset\left[V_{k+1, a}\right]$ for $a \in \mathcal{N}_{1}^{n} \cup \ldots \cup \mathcal{N}_{r}^{n}$.

We will often skip the arguments of $V_{k, a}$ when they are obvious from the context.
The values of $\left[x_{k+1}\right]$ and $\left[V_{k+1, a}\right], a \in \mathcal{N}_{1}^{n}$ are computed using one step of the $\mathcal{C}^{1}$-Lohner algorithm. After it is done, we perform the following operations to compute $\left[V_{k+1, a}\right]$ for $a \in \mathcal{N}_{2}^{n} \cup \ldots \cup \mathcal{N}_{r}^{n}$

1. Find a rough enclosure for $D_{a} \varphi\left(\left[0, h_{k}\right],\left[x_{k}\right]\right)$.
2. Compute $\left[V_{k+1, a}\right]$, this will also involve some rearrangement computations to reduce the wrapping effect for $V[\mathrm{Mo}, \mathrm{Lo}]$.

## 4 Computation of a rough enclosure for $D_{a} \varphi$

For a fixed multipointer $a \in \mathcal{N}_{d}^{n}$ Equation (16) can be written as follows

$$
\begin{equation*}
\frac{d}{d t} D_{a} \varphi(t, x)=B_{a}(t, x)+A(t, x) D_{a} \varphi(t, x) \tag{18}
\end{equation*}
$$

where

$$
\begin{align*}
B_{a} & =\sum_{k=2}^{d} \sum_{i_{1}, \ldots, i_{k}=1}^{n}\left(D^{e_{i_{1}}+\cdots+e_{i_{k}}} f\right) \circ \varphi \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{d}(k)} \prod_{j=1}^{k} D_{a_{\delta_{j}}} \varphi_{i_{j}},  \tag{19}\\
A & =D f \circ \varphi .
\end{align*}
$$

The procedure for computing the rough enclosure is based on the notion of the logarithmic norm.

Definition 9 [HNW] For a square matrix A the logarithmic norm $\mu(A)$ is defined as a limit

$$
\mu(A)=\limsup _{h \rightarrow 0^{+}} \frac{\|\mathrm{Id}+A h\|-1}{h},
$$

where $\|\cdot\|$ is a given matrix norm.
The formulas for the logarithmic norm of a real matrix in the most frequently used norms are (see [HNW])

1. for $\|x\|_{1}=\sum_{i}\left|x_{i}\right|, \mu(A)=\max _{j}\left(a_{j j}+\sum_{i \neq j}\left|a_{i j}\right|\right)$,
2. for $\|x\|_{2}=\sqrt{\sum_{i}\left|x_{i}\right|^{2}}, \mu(A)$ is equal to the largest eigenvalue of $(A+$ $\left.A^{T}\right) / 2$,
3. for $\|x\|_{\infty}=\max _{i}\left|x_{i}\right|, \mu(A)=\max _{i}\left(a_{i i}+\sum_{j \neq i}\left|a_{i j}\right|\right)$.

In order to find bounds for $D_{a} \varphi$ we use the following theorem [HNW, Thm. I.10.6]

Theorem 10 Let $x(t)$ be a solution of a differential equation

$$
\begin{equation*}
x^{\prime}(t)=f(t, x(t)), \quad x \in \mathbb{R}^{n} . \tag{20}
\end{equation*}
$$

Let $\nu(t)$ be a piecewise differentiable function with values in $\mathbb{R}^{n}$. Assume that

$$
\begin{array}{r}
\mu\left(\frac{\partial f}{\partial x}(t, \eta)\right) \leq l(t) \quad \text { for } \eta \in[x(t), \nu(t)] \\
\left|\nu^{\prime}(t)-f(t, \nu(t))\right| \leq \delta(t)
\end{array}
$$

Then for $t \geq t_{0}$ we have

$$
\begin{equation*}
|x(t)-\nu(t)| \leq e^{L(t)}\left(\left|x\left(t_{0}\right)-\nu\left(t_{0}\right)\right|+\int_{t_{0}}^{t} e^{-L(s)} \delta(s) d s\right) \tag{21}
\end{equation*}
$$

with $L(t)=\int_{t_{0}}^{t} l(\tau) d \tau$.
We apply the above theorem to Equation (18) to obtain

Lemma 11 Let us fix $x \in \mathbb{R}^{n}$. Assume that $\left|B_{a}(t, x)\right| \leq \delta(t)$ and $\mu(A(t, x)) \leq$ $l(t)$, then for $t>t_{0}$ holds

$$
\begin{equation*}
\left|D_{a} \varphi(t, x)\right| \leq\left|D_{a} \varphi\left(t_{0}, x\right)\right| e^{L(t)}+e^{L(t)} \int_{t_{0}}^{t} e^{-L(\tau)} \delta(\tau) d \tau \tag{22}
\end{equation*}
$$

with $L(t)=\int_{t_{0}}^{t} l(\tau) d \tau$.
Proof: Consider Equation (18) and a homogenous problem for (18)

$$
\begin{equation*}
\frac{d}{d t} w=f(t, w):=A(t, x) \cdot w, \quad w \in \mathbb{R}^{n} \tag{23}
\end{equation*}
$$

Using Theorem 10 we can estimate the difference between any solution of (23), $w$, and a solution of (18), denoted by $D_{a} \varphi$.

$$
\begin{equation*}
\left|D_{a} \varphi(t)-w(t)\right| \leq\left|D_{a} \varphi\left(t_{0}\right)-w\left(t_{0}\right)\right| e^{L(t)}+e^{L(t)} \int_{t_{0}}^{t} e^{-L(\tau)} \delta(\tau) d \tau \tag{24}
\end{equation*}
$$

After a substitution $w(t)=0$, which is a solution of the homogenous equation, we obtain our assertion.

Usually we do not have any control over the time dependence of $\delta$ and $l$, hence we will use the following

Lemma 12 Assume that $\left|B_{a}(t, x)\right| \leq \delta$ and $\mu(A(t, x)) \leq l$ for $t \in[0, h]$ then for $t \in[0, h]$ we have

$$
\begin{equation*}
\left|D_{a} \varphi(t, x)\right| \leq\left|D_{a} \varphi(0, x)\right| \max \left(1, e^{h l}\right)+\delta \frac{e^{l t}-1}{l}, \quad \text { if } l \neq 0 \tag{25}
\end{equation*}
$$

or

$$
\begin{equation*}
\left|D_{a} \varphi(t, x)\right| \leq\left|D_{a} \varphi(0, x)\right|+\delta t, \quad \text { when } l=0 \tag{26}
\end{equation*}
$$

### 4.1 The procedure for the computation of the rough enclosure for $V$.

For $a \in \mathcal{N}_{1}^{n} \cup \ldots \cup \mathcal{N}_{r}^{n}$ by $\left[E_{a}\right]$ we will denote the rough enclosure for the corresponding variational equation. The procedure for the computation of the rough enclosure $\left[E_{a}\right.$ ] is iterative, which means that given the rough enclosure for $\varphi\left(\left[0, h_{k}\right],\left[x_{k}\right]\right)$ and the rough enclosures $D_{a} \varphi\left(\left[0, h_{k}\right],\left[x_{k}\right]\right)$ for all $a \in \mathcal{N}_{1}^{n} \cup$ $\ldots \cup \mathcal{N}_{p}^{n}$ we are able to compute the rough enclosure for $D_{a} \varphi\left(\left[0, h_{k}\right],\left[x_{k}\right]\right)$ for $a \in \mathcal{N}_{p+1}^{n}$.

The procedures for computation of the rough enclosures of $\varphi\left(\left[0, h_{k}\right],\left[x_{k}\right]\right)$ and $D_{a} \varphi\left(\left[0, h_{k}\right],\left[x_{k}\right]\right)$ for $a \in \mathcal{N}_{1}^{n}$ have been given in [Z]. Below we present an algorithm for computing $\left[E_{a}\right]$ for $a \in \mathcal{N}_{2}^{n} \cup \ldots \cup \mathcal{N}_{r}^{n}$.
Input parameters:

- $h_{k}$ - a time step,
- $\left[x_{k}\right] \subset \mathbb{R}^{n}$ - the current value of $x=\varphi\left(t_{k},\left[x_{0}\right]\right)$,
- $\left[E_{0}\right] \subset \mathbb{R}^{n}$ - a compact and convex such that $\varphi\left(\left[0, h_{k}\right],\left[x_{k}\right]\right) \subset\left[E_{0}\right]$,
- $\left[E_{a}\right] \subset \mathbb{R}^{n}, a \in \mathcal{N}_{1}^{n} \cup \ldots \cup \mathcal{N}_{p}^{n}$ such that $D_{a} \varphi\left(\left[0, h_{k}\right],\left[x_{k}\right]\right) \subset\left[E_{a}\right]$ for $a \in \mathcal{N}_{1}^{n} \cup \ldots \cup \mathcal{N}_{p}^{n}$.


## Output:

- $\left[E_{a}\right] \subset \mathbb{R}^{n}, a \in \mathcal{N}_{p+1}^{n}$ such that

$$
D_{a} \varphi\left(\left[0, h_{k}\right],\left[x_{k}\right]\right) \subset\left[E_{a}\right]
$$

Before we present an algorithm let us observe that for a fixed $a \in \mathcal{N}_{p+1}^{n}, B_{a}$ defined in (19) could be seen as a multivariate function of $t, x$ and $V_{b}=D_{b} \varphi$ for $b \in \mathcal{N}_{1}^{n} \cup \ldots \cup \mathcal{N}_{p}^{n}$. More precisely, put $m_{p}:=\sharp\left\{\mathcal{N}_{1}^{n} \cup \ldots \cup \mathcal{N}_{p}^{n}\right\}$, where $\sharp$ stands for the number of elements of a set. Recall that we have defined by (7) a linear order in $\mathcal{N}^{n}$. Hence, there is a unique sequence of multipointers $b_{1}, \ldots, b_{m_{p}}$, such that $b_{i} \in \mathcal{N}_{1}^{n} \cup \ldots \cup \mathcal{N}_{p}^{n}$ for $i=1, \ldots, m_{p}, b_{1} \leq b_{2} \leq \cdots \leq b_{m_{p}}$ and $b_{i} \neq b_{j}$ for $i \neq j$.

Let us define

$$
\begin{aligned}
\tilde{B}_{a} & : \mathbb{R} \times\left(\mathbb{R}^{n}\right)^{m_{p}+1} \rightarrow \mathbb{R}^{n} \\
F_{a} & : \mathbb{R} \times\left(\mathbb{R}^{n}\right)^{m_{p}+1} \rightarrow \mathbb{R}^{n}
\end{aligned}
$$

by

$$
\begin{align*}
& \tilde{B}_{a}\left(t, x, v_{b_{1}}, \ldots, v_{b_{m_{p}}}\right)= \\
& \quad \sum_{k=2}^{p+1} \sum_{i_{1}, \ldots, i_{k}=1}^{n} D^{e_{i_{1}}+\cdots+e_{i_{k}}} f(\varphi(t, x)) \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p+1}(k)} \prod_{j=1}^{k}\left(v_{a_{\delta_{j}}}\right)_{i_{j}} \tag{27}
\end{align*}
$$

and

$$
\begin{equation*}
F_{a}\left(t, x, v_{b_{1}}, \ldots, v_{b_{m}}\right)=\tilde{B}_{a}\left(t, x, v_{b_{1}}, \ldots, v_{b_{m}}\right)+D f(\varphi(t, x)) V_{a}(t, x) \tag{28}
\end{equation*}
$$

## Algorithm:

To compute $\left[E_{a}\right]$ for $a \in \mathcal{N}_{p+1}^{n}$ we proceed as follows

1. Find $l \geq\left(\max _{x \in\left[E_{0}\right]} \mu(D f(x))\right)$.
2. Compute $\delta_{a} \geq \max \left\|\tilde{B}_{a}\right\|$, i.e.

$$
\delta_{a} \geq \max _{\left(x, v_{b_{1}}, \ldots, v_{b_{m_{p}}}\right) \in\left[E_{0}\right] \times\left[E_{b_{1}}\right] \times \cdots \times\left[E_{b_{m_{p}}}\right]}\left\|\tilde{B}_{a}\left(0, x, v_{b_{1}}, \ldots, v_{b_{m_{p}}}\right)\right\|
$$

For example, if $a=(j, c) \in \mathcal{N}_{2}^{n}$, then $\delta_{a}$ should be such that

$$
\delta_{a} \geq \max _{x \in\left[E_{0}\right], v_{1} \in\left[E_{(1)}\right], \ldots, v_{n} \in\left[E_{(n)}\right]}\left\|\sum_{r, s=1}^{n} \frac{\partial^{2} f}{\partial x_{r} \partial x_{s}}(x)\left(v_{j}\right)_{s}\left(v_{c}\right)_{r}\right\|
$$

3. Define $\left[E_{a}\right]_{i}=[-1,1] \delta_{a} \frac{e^{l t}-1}{l}$, for $i=1, \ldots, n$, where $\left[E_{a}\right]_{i}$ denotes $i$-th coordinate of $\left[E_{a}\right]$.

One can refine the obtained enclosure by

$$
\left[E_{a}\right]:=\left(\left[0, h_{k}\right] F_{a}\left(0,\left[E_{0}\right],\left[E_{b_{1}}\right], \ldots,\left[E_{b_{m_{p}}}\right]\right)\right) \cap\left[E_{a}\right] .
$$

Indeed, from (17) we have $D_{a} \varphi_{i}\left(0, x_{0}\right)=0$ for $i=1, \ldots, n, x_{0} \in\left[E_{0}\right]$ and for $t \in\left[0, h_{k}\right]$ we have

$$
\begin{aligned}
D_{a} \varphi_{i}\left(t, x_{0}\right) & =D_{a} \varphi_{i}\left(t, x_{0}\right)-D_{a} \varphi_{i}\left(0, x_{0}\right) \\
& =t\left(F_{a}\right)_{i}\left(\theta_{i}, x_{0}, D_{b_{1}} \varphi\left(\theta_{i}, x_{0}\right), \ldots, D_{b_{m_{p}}} \varphi\left(\theta_{i}, x_{0}\right)\right) \\
& =t\left(F_{a}\right)_{i}\left(0, \varphi\left(\theta_{i}, x_{0}\right), D_{b_{1}} \varphi\left(\theta_{i}, x_{0}\right), \ldots, D_{b_{m_{p}}} \varphi\left(\theta_{i}, x_{0}\right)\right)
\end{aligned}
$$

for some $\theta_{i} \in[0, t] \subset\left[0, h_{k}\right]$. In the above we have used the fact that

$$
F_{a}\left(t, x, v_{1}, \ldots, v_{m_{p}}\right)=F_{a}\left(0, \varphi(t, x), v_{1}, \ldots, v_{m_{p}}\right)
$$

Since $\varphi\left(\theta_{i}, x_{0}\right) \in\left[E_{0}\right]$ and $D_{b_{j}} \varphi\left(\theta_{i}, x_{0}\right) \in\left[E_{b_{j}}\right]$ for $j=1, \ldots, m_{p}$ we get

$$
D_{a} \varphi_{i}\left(t, x_{0}\right) \in\left[0, h_{k}\right]\left(F_{a}\right)_{i}\left(0,\left[E_{0}\right],\left[E_{b_{1}}\right], \ldots,\left[E_{b_{m_{p}}}\right]\right)
$$

## 5 Computation of [ $V_{k+1}$ ]

### 5.1 Composition formulas

We apply the Faá di Bruno formula (10) to $f=\varphi\left(h_{k}, \cdot\right)$ and $g=\varphi\left(t_{k}, \cdot\right)$ to obtain

$$
\begin{array}{r}
V_{a}\left(t_{k}+h_{k}, x_{0}\right)=\sum_{k=1}^{p} \sum_{i_{1}, \ldots, i_{k}=1}^{n} V_{\Lambda^{-1}\left(e_{i_{1}}+\ldots+e_{i_{k}}\right)}\left(h_{k}, x_{k}\right) \\
\sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p}(k)} \prod_{j=1}^{k}\left(V_{a_{\delta_{j}}}\right)_{i_{j}}\left(t_{k}, x_{0}\right)
\end{array}
$$

for all $x_{0} \in\left[x_{0}\right]$. Using notations $\left[V_{k+1, a}\right]:=\left[V_{a}\left(t_{k}+h_{k},\left[x_{0}\right]\right)\right]$ and $\left[V_{k, a}\right]=$ [ $\left.V_{a}\left(t_{k},\left[x_{0}\right]\right)\right]$ we can rewrite the above equation as

$$
\begin{equation*}
\left[V_{k+1, a}\right]=\sum_{k=1}^{p} \sum_{i_{1}, \ldots, i_{k}=1}^{n} V_{\Lambda^{-1}\left(e_{i_{1}}+\ldots+e_{i_{k}}\right)}\left(h_{k},\left[x_{k}\right]\right) \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p}(k)} \prod_{j=1}^{k}\left[V_{k, a_{\delta_{j}}}\right]_{i_{j}} \tag{29}
\end{equation*}
$$

where $\Lambda$ is defined by (6).

### 5.2 The procedure for computation of [ $V_{k+1}$ ]

We introduce new parameters $o_{d}$ - the order of the Taylor method used in computations of $V_{a}$ for $a \in \mathcal{N}_{d}^{n}$. It makes sense to take $o_{1} \geq o_{2} \geq \cdots \geq o_{r}$.
Input parameters:

- $h_{k}$ - a time step,
- $\left[x_{k}\right] \subset \mathbb{R}^{n}$ - the current value of $x=\varphi\left(t_{k},\left[x_{0}\right]\right)$,
- $\left[V_{k, a}\right] \subset \mathbb{R}^{n}$ - a current value of $V_{k, a}\left(t_{k},\left[x_{0}\right]\right)$, for $a \in \mathcal{N}_{1}^{n} \cup \ldots \cup \mathcal{N}_{r}^{n}$,
- $\left[E_{0}\right] \subset \mathbb{R}^{n}$ compact and convex, such that $\varphi\left(\left[0, h_{k}\right],\left[x_{k}\right]\right) \subset\left[E_{0}\right]$ - a rough enclosure for $\left[x_{k}\right]$,
- $\left[E_{a}\right] \subset \mathbb{R}^{n}$, compact and convex, such that $D_{a} \varphi\left(\left[0, h_{k}\right],\left[x_{k}\right]\right) \subset\left[E_{a}\right]$, for $a \in \mathcal{N}_{1}^{n} \cup \ldots \cup \mathcal{N}_{r}^{n}$.

Output: $\left[V_{k+1, a}\right] \subset \mathbb{R}^{n}$, such that

$$
\begin{equation*}
V_{a}\left(t_{k}+h_{k}, x_{0}\right) \in\left[V_{k+1, a}\right] \tag{30}
\end{equation*}
$$

for $x_{0} \in\left[x_{0}\right]$ and $a \in \mathcal{N}_{1}^{n} \cup \ldots \cup \mathcal{N}_{r}^{n}$.
Algorithm: We compute [ $V_{k+1}$ ] as follows

1. Computation of $V_{a}\left(h_{k},\left[x_{k}\right]\right)$ using the Taylor method for Equation (16), i.e. for $a \in \mathcal{N}_{p}^{n}$ we compute

$$
\begin{align*}
{\left[F_{a}\right] } & =\sum_{i=1}^{o_{p}} \frac{h_{k}^{i}}{i!} \frac{d^{i-1}}{d t^{i-1}} F_{a}\left(0,\left[x_{k}\right], V_{b_{1}}, \ldots, V_{b_{m_{p-1}}}\right)  \tag{31}\\
& +\frac{h^{o_{p}+1}}{\left(o_{p}+1\right)!} \frac{d^{o_{p}}}{d t^{o_{p}}} F_{a}\left(0,\left[E_{0}\right],\left[E_{b_{1}}\right], \ldots,\left[E_{b_{m_{p-1}}}\right]\right)
\end{align*}
$$

where $V_{b_{i}}=0$ for $b_{i} \in \mathcal{N}_{2}^{n} \cup \ldots \cup \mathcal{N}_{p-1}^{n}$ and $V_{(j)}=e_{j}^{n}$ for $j=1, \ldots, n$. Observe that

$$
\begin{equation*}
V_{a}\left(h_{k},\left[x_{k}\right]\right) \subset\left[F_{a}\right] . \tag{32}
\end{equation*}
$$

Indeed, using the Taylor series expansion we obtain that for $x_{k} \in\left[x_{k}\right]$ and $j=1, \ldots, n$ holds

$$
\begin{array}{r}
\left(V_{a}\right)_{j}\left(h_{k}, x_{k}\right)=\sum_{i=1}^{o_{p}} \frac{h_{k}^{i}}{i!} \frac{d^{i-1}}{d t^{i-1}}\left(F_{a}\right)_{j}\left(0, x_{k}, V_{b_{1}}\left(0, x_{k}\right), \ldots, V_{b_{m_{p-1}}}\left(0, x_{k}\right)\right) \\
\quad+\frac{h^{o_{p}+1}}{\left(o_{p}+1\right)!} \frac{d^{o_{p}}}{d t^{o_{p}}}\left(F_{a}\right)_{j}\left(\theta_{i}, x_{k}, V_{b_{1}}\left(\theta_{i}, x_{k}\right), \ldots, V_{b_{m_{p-1}}}\left(\theta_{i}, x_{k}\right)\right)
\end{array}
$$

for some $\theta_{i} \in\left[0, h_{k}\right]$. Observe, that

$$
\begin{aligned}
& \frac{d^{o_{p}}}{d t^{o_{p}}}\left(F_{a}\right)_{j}\left(\theta_{i}, x_{k}, V_{b_{1}}\left(\theta_{i}, x_{k}\right), \ldots, V_{b_{m_{p-1}}}\left(\theta_{i}, x_{k}\right)\right)= \\
& \quad \frac{d^{o_{p}}}{d t^{o_{p}}}\left(F_{a}\right)_{j}\left(0, \varphi\left(\theta_{i}, x_{k}\right), V_{b_{1}}\left(\theta_{i}, x_{k}\right), \ldots, 0, V_{b_{m_{p-1}}}\left(\theta_{i}, x_{k}\right)\right) .
\end{aligned}
$$

Since $\varphi\left(\theta_{i}, x_{k}\right) \in\left[E_{0}\right]$ and $V_{b_{s}}\left(\theta_{i}, x_{k}\right) \in\left[E_{b_{s}}\right]$ for $s=1, \ldots, m_{p-1}$ we obtain our assertion.
2. The composition. Put

$$
\left[J_{k}\right]:=\left(\left[F_{(1)}\right], \ldots,\left[F_{(n)}\right]\right)^{T} .
$$

Using (29) for $a \in \mathcal{N}_{p}^{n}$ we have

$$
\begin{equation*}
\left[V_{k+1, a}\right]=\left[\alpha_{a}\right]+\left[J_{k}\right] \cdot\left[V_{k, a}\right], \tag{33}
\end{equation*}
$$

where

$$
\begin{equation*}
\left[\alpha_{a}\right]=\sum_{k=2}^{p} \sum_{i_{1}, \ldots, i_{k}=1}^{n}\left[F_{\left.\Lambda^{-1}\left(e_{i_{1}}+\ldots+e_{i_{k}}\right)\right]} \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p}(k)} \prod_{j=1}^{k}\left[V_{k, a_{\delta_{j}}}\right]_{i_{j}}\right. \tag{34}
\end{equation*}
$$

### 5.3 Rearrangement for $V_{a}$ - the evaluation of Equation (33)

It is well known that a direct evaluation of Equation (33) leads to the wrapping effect [Mo, Lo]. To avoid it, following the work of Lohner [Lo] we will use the scheme proposed in $[\mathrm{Z}]$ for the $\mathcal{C}^{1}$-algorithm.

Namely, observe that Equation (33) has exactly the same structure as the propagation equations for $\mathcal{C}^{1}$-method (see [Z, Section 3]). Moreover, all vectors $V_{k, a}$, for $a \in \mathcal{N}_{1}^{n} \cup \ldots \mathcal{N}_{r}^{n}$ 'propagate' by the same $\left[J_{k}\right]$ as does the variational part in [Z], hence it makes sense to use the same approach.

To be more precise, each set $\left[V_{k, a}\right]$, for $a \in \mathcal{N}_{1}^{n} \cup \ldots \cup \mathcal{N}_{r}^{n}$, is represented in the following form

$$
\left[V_{k, a}\right]=v_{k, a}+\left[B_{k}\right]\left[r_{k, a}\right]+C_{k}\left[q_{k, a}\right],
$$

where $\left[B_{k}\right]$ is an interval matrix, $C_{k}$ is a point matrix, $v_{k, a}$ is a point vector and $r_{k, a}, q_{k, a}$ are interval vectors. Observe that $\left[B_{k}\right]$ and $C_{k}$ are independent of $a$.

In the sequel we will drop the index $a$. Equation (33) leads to

$$
\begin{equation*}
\left[V_{k+1}\right]=[\alpha]+\left[J_{k}\right]\left(v_{k}+\left[B_{k}\right]\left[r_{k}\right]+C_{k}\left[q_{k}\right]\right) \tag{35}
\end{equation*}
$$

Let $m([z])$ denotes a center of an interval object, i.e. $[z]$ is interval vector or interval matrix and let $\Delta([z])=[z]-m([z])$.

Let $[Q]$ be an interval matrix which contains an orthogonal matrix. Usually, $[Q]$ is computed by the orthonormalisation of the columns of $m\left(\left[J_{k}\right]\left[B_{k}\right]\right)$.

Let

$$
\begin{aligned}
{[Z] } & =\left[J_{k}\right] C_{k}, \\
C_{k+1} & =m([Z]), \\
{\left[B_{k+1}\right] } & =[Q] .
\end{aligned}
$$

Then we rearrange formula (35) as follows

$$
\begin{align*}
{[s] } & =[\alpha]+\left[J_{k}\right] v_{k}+\Delta([Z])\left[q_{k}\right], \\
v_{k+1} & =m([s]),  \tag{36}\\
{\left[q_{k+1}\right] } & =\left[q_{k}\right] \\
{\left[r_{k+1}\right] } & =\left[Q^{T}\right] \Delta([s])+\left(\left[Q^{T}\right]\left[J_{k}\right]\left[B_{k}\right]\right)\left[r_{k}\right] .
\end{align*}
$$

Summarizing, we can use the following data structure to represent $\varphi\left(t_{k},\left[x_{0}\right]\right)$ and $D_{a} \varphi\left(t_{k},\left[x_{0}\right]\right)$, for $a \in \mathcal{N}_{1}^{n} \cup \ldots \cup \mathcal{N}_{r}^{n}$

```
type \(\mathrm{CnSet}=\) record
    \(v_{0}, r_{0}, q_{0}\) : IntervalVector;
    \(C_{0}, B_{0}, C, B\) : IntervalMatrix;
    \(\left\{v_{a}, r_{a}, q_{a} \text { : IntervalVector }\right\}_{a \in \mathcal{N}_{1}^{n} \cup \ldots \cup \mathcal{N}_{r}^{n}}\)
```

end;

The set $\varphi\left(t_{k},\left[x_{0}\right]\right)$ is represented as $v_{0}+B_{0} r_{0}+C_{0} q_{0}$, the partial derivatives $D_{a} \varphi\left(t_{k},\left[x_{0}\right]\right)$ are represented as $v_{a}+B r_{a}+C q_{a}$. The matrices $B, C$ are common for all partial derivatives.

Notice, that if we start the $\mathcal{C}^{r}$ computation with an initial condition (17) then there is no Lipschitz part at the beginning for the partial derivatives. Hence, the initial values for $C$ and $B$ are set to the identity matrix and the initial values for $q_{a}, r_{a}$ are set to zero.

If the interval vectors $r_{a}$ become 'thick' (i.e. theirs diameters are larger than some threshold value) we can set a new Lipschitz part in our representation (it must be done simultaneously for all $D_{a} \varphi$ ) and reset $r_{a}$ in the following way

$$
\begin{aligned}
q_{a} & =r_{a}+\left(B^{T} C\right) q_{a}, \quad \text { for } a \in \mathcal{N}_{1}^{n} \cup \ldots \cup \mathcal{N}_{r}^{n} \\
r_{a} & =0, \quad \text { for } a \in \mathcal{N}_{1}^{n} \cup \ldots \cup \mathcal{N}_{r}^{n} \\
C & =B, \\
B & =\text { Id. }
\end{aligned}
$$

This is a place where a discontinuity (non-monotonicity) appears in the algorithm. A similar change of the Lipschitz part may be done when vectors $r_{a}$ become thick in comparison to $q_{a}$.

## 6 Derivatives of Poincaré map

Consider a differential equation

$$
\begin{equation*}
x^{\prime}=f(x), \quad x \in \mathbb{R}^{n}, \quad f \in \mathcal{C}^{K+1} \tag{37}
\end{equation*}
$$

Let $\varphi: \mathbb{R} \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ be a (local) dynamical system induced by (37). Let $\alpha: \mathbb{R}^{n} \rightarrow \mathbb{R}$ be $\mathcal{C}^{1}$-map. Put $\Pi=\{x \mid \alpha(x)=C\}$.

Definition 13 We will say that $\Pi$ is a local section for the vector field $f$ at $y_{0} \in \Pi$ if

$$
\begin{equation*}
\left\langle\nabla \alpha\left(y_{0}\right) \mid f\left(y_{0}\right)\right\rangle \neq 0 \tag{38}
\end{equation*}
$$

Assume $x_{0} \in \mathbb{R}^{n}$ and $t_{0} \in \mathbb{R}$ are such that $\Pi$ is a local section at $\varphi\left(t_{0}, x_{0}\right)$. Consider an implicit equation

$$
\begin{equation*}
\alpha\left(\varphi\left(t_{P}(x), x\right)\right)=C \tag{39}
\end{equation*}
$$

It follows easily from (38) and from the implicit function theorem that there exists a uniquely defined $t_{P}: \mathbb{R}^{n} \rightarrow \rightarrow \mathbb{R}$ in a neighborhood of $x_{0}$, such that $t_{P}\left(x_{0}\right)=t_{0}$. The function $t_{P}$ is as smooth as the flow $\varphi$. We will refer to $t_{P}$ as to the Poincare return time to section $\Pi$.

We define a Poincaré map $P: \mathbb{R}^{n} \supset \operatorname{dom}\left(t_{P}\right) \rightarrow \mathbb{R}^{n}$ by

$$
\begin{equation*}
P(x)=\varphi\left(t_{P}(x), x\right) . \tag{40}
\end{equation*}
$$

Usually the Poincaré map is defined as a map $P: \Pi_{1} \rightarrow \Pi_{2}$, where $\Pi_{1}, \Pi_{2}$ are some local sections in $\mathbb{R}^{n}$. The approach taken here, i.e. treating the Poincaré map as map $P: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ allows us to not to worry about the coordinates on the local section.

We are interested in the partial derivatives of $P$ defined by (40).
From (40) we obtain

$$
\begin{equation*}
\frac{\partial P_{i}}{\partial x_{j}}(x)=f_{i}(P(x)) \frac{\partial t_{P}}{\partial x_{j}}(x)+\frac{\partial \varphi_{i}}{\partial x_{j}}\left(t_{P}(x), x\right) \tag{41}
\end{equation*}
$$

We need $\frac{\partial t_{P}}{\partial x_{j}}$. We differentiate (39) to obtain

$$
\begin{align*}
\sum_{k=1}^{n} \frac{\partial \alpha}{\partial x_{k}}(P(x))\left(f_{k}(P(x)) \frac{\partial t_{P}}{\partial x_{j}}(x)+\frac{\partial \varphi_{k}}{\partial x_{j}}\left(t_{P}(x), x\right)\right) & =0, \\
\langle\nabla \alpha(P(x)) \mid f(P(x))\rangle \frac{\partial t_{P}}{\partial x_{j}}(x)+\sum_{k=1}^{n} \frac{\partial \alpha}{\partial x_{k}}(P(x)) \frac{\partial \varphi_{k}}{\partial x_{j}}\left(t_{P}(x), x\right) & =0 . \tag{42}
\end{align*}
$$

Hence

$$
\begin{equation*}
\frac{\partial t_{P}}{\partial x_{j}}(x)=-\frac{1}{\langle\nabla \alpha(P(x)) \mid f(P(x))\rangle} \sum_{k=1}^{n} \frac{\partial \alpha}{\partial x_{k}}(P(x)) \frac{\partial \varphi_{k}}{\partial x_{j}}\left(t_{P}(x), x\right) . \tag{43}
\end{equation*}
$$

### 6.1 Higher order derivatives of the Poincaré map

To make formulas transparent we will drop arguments of functions in this section, but reader should be aware that for $t_{P}$ and its partial derivatives the argument is $x$, for $\varphi$ and $D_{a} \varphi$ the argument is always the pair $\left(t_{P}(x), x\right)$.

From (41) we obtain

$$
\begin{aligned}
D_{(j, c)} P & =\frac{\partial^{2}}{\partial t^{2}} \varphi D_{(j)} t_{P} D_{(c)} t_{P}+\frac{\partial}{\partial t} D_{(c)} \varphi D_{(j)} t_{P}+\frac{\partial}{\partial t} \varphi D_{(j, c)} t_{P} \\
& +\frac{\partial}{\partial t} D_{(j)} \varphi D_{(c)} t_{P}+D_{(j, c)} \varphi
\end{aligned}
$$

It is easy to see that partial derivatives of high order give rise to quite complex expressions and it is not entirely obvious how to organize it in some coherent and programmable way. For this purpose we use the following

Lemma 14 For a multipointer $a \in \mathcal{N}_{p}^{n}$ we have

$$
\begin{align*}
D_{a} P & =D_{a} \varphi+\frac{\partial \varphi}{\partial t} D_{a} t_{P} \\
& +\sum_{k=2}^{p} \frac{\partial^{k} \varphi}{\partial t^{k}} \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p}(k)} \prod_{j=1}^{k} D_{a_{\delta_{j}}} t_{P}  \tag{44}\\
& +\sum_{k=2}^{p} \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p}(k)} \sum_{s=1}^{k} \frac{\partial^{k-1}}{\partial t^{k-1}} D_{a_{\delta_{s}}} \varphi \prod_{j \neq s} D_{a_{\delta_{j}}} t_{P} .
\end{align*}
$$

Proof: By induction on $p$. For $p=1$ formula (44) is equivalent to (41), because the two last sums are taken over empty set. Assume (44) holds true for some $p \geq 1$ and fix $a \in \mathcal{N}_{p+1}^{n}$. Our goal is to show that

$$
D_{a} P=R_{1}+R_{2}+R_{3}
$$

where

$$
\begin{aligned}
R_{1} & =D_{a} \varphi+\frac{\partial}{\partial t} \varphi D_{a} t_{P} \\
R_{2} & =\sum_{k=2}^{p+1} \frac{\partial^{k}}{\partial t^{k}} \varphi \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p+1}(k)} \prod_{j=1}^{k} D_{a_{\delta_{j}}} t_{P} \\
R_{3} & =\sum_{k=2}^{p+1} \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p+1}(k)} \sum_{s=1}^{k} \frac{\partial^{k-1}}{\partial t^{k-1}} D_{a_{\delta_{s}}} \varphi \prod_{j \neq s} D_{a_{\delta_{j}}} t_{P} .
\end{aligned}
$$

Write $a=\beta+\gamma$, where $\beta \in \mathcal{N}_{p}^{n}$ and $\gamma=\left(a_{p+1}\right) \in \mathcal{N}_{1}^{n}$. From the induction assumption we have

$$
\begin{aligned}
D_{a} P & =D_{\gamma}\left(D_{\beta} \varphi+\frac{\partial}{\partial t} \varphi D_{\beta} t_{P}\right) \\
& +D_{\gamma}\left(\sum_{k=2}^{p} \frac{\partial^{k}}{\partial t^{k}} \varphi \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p}(k)} \prod_{j=1}^{k} D_{\beta_{\delta_{j}}} t_{P}\right) \\
& +D_{\gamma}\left(\sum_{k=2}^{p} \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p}(k)} \sum_{s=1}^{k} \frac{\partial^{k-1}}{\partial t^{k-1}} D_{\beta_{\delta_{s}}} \varphi \prod_{j \neq s} D_{\beta_{\delta_{j}}} t_{P}\right) \\
& =\sum_{i=1}^{10} S_{i},
\end{aligned}
$$

where

$$
\begin{aligned}
S_{1}= & D_{a} \varphi+\frac{\partial}{\partial t} \varphi D_{a} t_{P}, \\
S_{2}= & \frac{\partial}{\partial t} D_{\beta} \varphi D_{\gamma} t_{P}, \\
S_{3}= & \frac{\partial^{2}}{\partial t^{2}} \varphi D_{\beta} t_{P} D_{\gamma} t_{P}, \\
S_{4}= & \frac{\partial}{\partial t} D_{\gamma} \varphi D_{\beta} t_{P}, \\
S_{5}= & \sum_{k=2}^{p} \frac{\partial^{k}}{\partial t^{k}} D_{\gamma} \varphi \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p}(k)} \prod_{j=1}^{k} D_{\beta_{\delta_{j}}} t_{P}, \\
S_{6}= & \sum_{k=2}^{p} \frac{\partial^{k+1}}{\partial t^{k+1}} \varphi D_{\gamma} t_{P} \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p}(k)} \prod_{j=1}^{k} D_{\beta_{\delta_{j}}} t_{P}, \\
S_{7}= & \sum_{k=2}^{p} \frac{\partial^{k}}{\partial t^{k}} \varphi \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p}(k)} \sum_{s=1}^{k} D_{\beta_{\delta_{s}}+\gamma} t_{P} \prod_{\substack{j=1 \\
j \neq s}}^{k} D_{\beta_{\delta_{j}}} t_{P}, \\
S_{8}= & \sum_{k=2}^{p} \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p}(k)} \sum_{s=1}^{k} \frac{\partial^{k-1}}{\partial t^{k-1}} D_{\beta_{\delta_{s}}+\gamma} \varphi \prod_{j \neq s} D_{\beta_{\delta_{j}}} t_{P}, \\
S_{9}= & \sum_{k=2}^{p} \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p}(k)} \sum_{s=1}^{k} \frac{\partial^{k}}{\partial t^{k}} D_{\beta_{\delta_{s}}} \varphi D_{\gamma} t_{P} \prod_{j \neq s} D_{\beta_{\delta_{j}}} t_{P}, \\
S_{10}= & \sum_{k=2}^{p} \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p}(k)} \\
& \sum_{s=1}^{k} \sum_{\substack{r=1 \\
r \neq s}}^{k \partial^{k-1}} \partial t^{k-1} D_{\beta_{\delta_{s}}} \varphi D_{\beta_{\delta_{r}}+\gamma} t_{P} \prod_{\substack{j \neq s \\
j \neq r}} D_{\beta_{\delta_{j}}} t_{P} .
\end{aligned}
$$

Obviously $R_{1}=S_{1}$. We will show that $R_{2}=S_{3}+S_{6}+S_{7}$ and $R_{3}=S_{2}+S_{4}+$ $S_{5}+S_{8}+S_{9}+S_{10}$.

Denote by $R_{i, k}, i=2,3$ a part of sum $R_{i}$ with fixed $k=2, \ldots, p+1$. Similarly, let us denote by $S_{i, k}$ a part of $\operatorname{sum} S_{i}, i=5, \ldots, 10$, for $k=2, \ldots, p$.

Using decomposition of $\mathcal{N}^{p+1}(2)$ as in (9) we obtain that $R_{2,2}=S_{3}+S_{7,2}$. Similarly, using (9) we observe that $R_{2, k}=S_{6, k-1}+S_{7, k}$ for $k=3, \ldots, p$. Finally, since $\mathcal{N}^{p+1}(p+1)=\{((1),(2), \ldots,(p+1))\}$ and $\gamma=\left(a_{p+1}\right)$ we find that $R_{2, p+1}=S_{6, p}$. This shows that $R_{2}=S_{3}+S_{6}+S_{7}$.

It remains to show that $R_{3}=S_{2}+S_{4}+S_{5}+S_{8}+S_{9}+S_{10}$. We will classify possible terms by the fact, where $p+1$ appears in $\delta_{i}, i=1, \ldots, k$ and how this $\delta_{i}$ enters in $R_{3}$ as $\delta_{s}$ or $\delta_{j}$. There are four cases

1. $\delta_{s}=(p+1)$,
2. $\delta_{j}=(p+1)$,
3. $p+1 \in \delta_{s},\left|\delta_{s}\right| \geq 2$,
4. $p+1 \in \delta_{j},\left|\delta_{j}\right| \geq 2$.

Let us fix $k=2$. Let $\left(\delta_{1}, \delta_{2}\right) \in \mathcal{N}^{p+1}(2)$. The term for case 1 is $S_{4}$, for case 2 is $S_{2}$, case 3 is $S_{8,2}$ and case 4 is $S_{10,2}$. Hence, $R_{3,2}=S_{2}+S_{4}+S_{8,2}+S_{10,2}$.

For $k=3, \ldots, p$ and fixed $\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p+1}(k)$ we have: case 1 is given by $S_{5, k-1}$, case 2 by $S_{9, k-1}$, case 3 by $S_{8, k}$ and case 4 by $S_{10, k}$ Hence, for $k=3, \ldots, p$ we have $R_{3, k}=S_{5, k-1}+S_{9, k-1}+S_{8, k}+S_{10, k}$.

Finally, for $k=p+1$ we observe, that $R_{3, p+1}=S_{5, p}+S_{9, p}$. Indeed, in this case $\left(\delta_{1}, \ldots, \delta_{p+1}\right)=((1),(2), \ldots,(p+1))$. Hence, either for $\delta_{s}=\gamma$ we have term $S_{5, p}$ and $\delta_{s} \neq \gamma$ we have $S_{9, p}$.

We have showed that $R_{3}=S_{2}+S_{4}+S_{5}+S_{8}+S_{9}+S_{10}$ and the proof is finished.

Hence, if we know all the partial derivatives of $t_{P}$ up order $p$ we can compute the partial derivatives of the Poincare map up the same order. In the next subsection we show how to compute partial derivatives of $t_{P}$ for affine sections.

### 6.2 Partial derivatives of $t_{P}$ for affine sections

Assume $\alpha: \mathbb{R}^{n} \rightarrow \mathbb{R}$ is an affine map given by

$$
\alpha(x)=\alpha_{0}+\sum_{i=1}^{n} \alpha_{i} x_{i} .
$$

This is a quite restrictive assumption about sections, but it leads to relatively simple formulas for $D_{a} t_{P}$ and it is sufficient for the applications we have in mind.

Lemma 15 For a multipointer $a \in \mathcal{N}_{p}^{n}$ holds

$$
\begin{aligned}
-D_{a} t_{P}\left\langle\nabla \alpha \left\lvert\, \frac{\partial}{\partial t} \varphi\right.\right\rangle & = \\
\left\langle\nabla \alpha \mid D_{a} \varphi\right\rangle & +\sum_{k=2}^{p}\left\langle\nabla \alpha \left\lvert\, \frac{\partial^{k}}{\partial t^{k}} \varphi\right.\right\rangle \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p}(k)} \prod_{j=1}^{k} D_{a_{\delta_{j}}} t_{P} \\
& +\sum_{k=2}^{p} \sum_{\left(\delta_{1}, \ldots, \delta_{k}\right) \in \mathcal{N}^{p}(k)} \sum_{s=1}^{k}\left\langle\nabla \alpha \left\lvert\, \frac{\partial^{k-1}}{\partial t^{k-1}} D_{a_{\delta_{s}}} \varphi\right.\right\rangle \prod_{j \neq s} D_{a_{\delta_{j}}} t_{P}
\end{aligned}
$$

Proof: The proof is a direct consequence of Lemma 14 and (39). Since $\alpha$ is affine, by differentiating of $\alpha(P(x))=C$ we get $\left\langle\nabla \alpha \mid D_{a} P\right\rangle=0$. Using formula (44) for $D_{a} P$ we obtain our assertion.

Fix $[x] \subset \mathbb{R}^{n}$ and assume we have a rigorous bound for $t_{P}([x]) \in\left[t_{1}, t_{2}\right]$ (see [Z, Section 6] for more details on this). Lemmas 15 and 14 show that given rigorous bounds for the partial derivatives $D_{a} \varphi\left(\left[t_{1}, t_{2}\right],[x]\right)$ and $\frac{\partial^{k}}{\partial t^{k}} D_{a} \varphi\left(\left[t_{1}, t_{2}\right],[x]\right)$ up to some order $p$ we can compute recursively rigorous bounds for the partial derivatives of $t_{P}([x])$ and $P([x])$ up to the same order. Notice, that $\frac{\partial^{k}}{\partial t^{k}} D_{a} \varphi$ are given by Taylor coefficients of the solution of (16) with initial conditions $P([x])$ for $\mathcal{C}^{0}$ part and $D_{a} \varphi\left(t_{P}(x),[x]\right)$ for equations for variations. Hence, these coefficients can be easily computed using the automatic differentiation algorithm.

## 7 Comparison to $\mathcal{C}^{0}$-solver.

In this section we present results of comparison of the $\mathcal{C}^{0}$-solver applied to the second order variational equations with $\mathcal{C}^{2}$-solver. We performed tests of these algorithms on some classical low dimensional examples, such as the VolterraLotka system

$$
\left\{\begin{array}{l}
\dot{x}=x(2-y),  \tag{45}\\
\dot{y}=y(x-3),
\end{array}\right.
$$

the pendulum equation

$$
\begin{equation*}
\ddot{x}=-\sin (x), \tag{46}
\end{equation*}
$$

the Lorenz system

$$
\left\{\begin{array}{l}
\dot{x}=10(-x+y),  \tag{47}\\
\dot{y}=28 x-y-x z, \\
\dot{z}=x y-\frac{8}{3} z,
\end{array}\right.
$$

the Michelson system

$$
\begin{equation*}
\dddot{x}+\dot{x}+\frac{1}{2} x^{2}=1 \tag{48}
\end{equation*}
$$

the Rössler system

$$
\left\{\begin{array}{l}
\dot{x}=-(y+z),  \tag{49}\\
\dot{y}=x+0.2 y, \\
\dot{z}=0.2+z(x-5.7),
\end{array}\right.
$$

and for the Hénon-Heiles system (Hamiltonian equation)

$$
\left\{\begin{array}{l}
\ddot{x}=-x-2 x y,  \tag{50}\\
\ddot{y}=y^{2}-y-x^{2} .
\end{array}\right.
$$

## General settings of the tests.

- We integrate the above systems together with second and third order variational equations along periodic orbits using $\mathcal{C}^{2}, \mathcal{C}^{3}$ and $\mathcal{C}^{0}$ solvers from the CAPD library [CAPD] to obtain bounds for the higher order derivatives. These periodic orbits are presented in Fig. 1. In each case the time of integration is equal to an approximate period of the orbit. We believe that this is a relevant time scale for the computer assisted proofs for these systems.
- When integrating the systems of variational equations using the $\mathcal{C}^{0}$ solver we simply add the variational equations to the main equations and apply the $\mathcal{C}^{0}$ solver to the extended system that has dimension $n\binom{n+k}{k}$, where $n$ is the dimension of the main problem and $k$ is order of derivatives we require.
- For each ODE (45)-(50) we set as initial conditions to each routine three boxes of diameters $0,10^{-10}$ and $10^{-6}$ centered at a point very close to the corresponding periodic solution. The actual initial conditions are given in the caption of Fig. 1.
- In each case we use the Taylor method of the order 20 with variable time step. The minimal acceptable time step has been set to $10^{-5}$. The computations were performed using the interval arithmetic with double precision.


Figure 1: Periodic orbits for the systems (45-50). The initial conditions are: $(2.5,1.5)$ for the Volterra-Lotka system, $(0.5,0.5)$ for pendulum equation, $(0,1.52596,0)$ for the Michelson system, $(-2.14737,2.07805,27)$ for the Lorenz system, ( $0 .,-8.3809417428298,0.029590060630665$ ) for the Rössler system and $(x, y, \dot{x}, \dot{y})=(0,0.10903,0,0.567723)$ for the Hénon-Heiles system.


Figure 2: Plots of $t \rightarrow r\left(\left[D^{2} \phi\left(x_{0}, t\right)\right]\right)$ and $t \rightarrow r\left(\left[D^{3} \phi\left(x_{0}, t\right)\right]\right)$ for the VolterraLotka system (45) obtained from $\mathcal{C}^{0}$ and $\mathcal{C}^{r}$ solvers for various diameters of initial conditions.


Figure 3: Plots of $t \rightarrow r\left(\left[D^{2} \phi\left(x_{0}, t\right)\right]\right)$ and $t \rightarrow r\left(\left[D^{3} \phi\left(x_{0}, t\right)\right]\right)$ for the pendulum equation (46) obtained from $\mathcal{C}^{0}$ and $\mathcal{C}^{r}$ solvers for various diameters of initial conditions.


Figure 4: Plots of $t \rightarrow r\left(\left[D^{2} \phi\left(x_{0}, t\right)\right]\right)$ and $t \rightarrow r\left(\left[D^{3} \phi\left(x_{0}, t\right)\right]\right)$ for the Michelson system (48) obtained from $\mathcal{C}^{0}$ and $\mathcal{C}^{r}$ solvers for various diameters of initial conditions.


Figure 5: Plots of $t \rightarrow r\left(\left[D^{2} \phi\left(x_{0}, t\right)\right]\right)$ and $t \rightarrow r\left(\left[D^{3} \phi\left(x_{0}, t\right)\right]\right)$ for the Lorenz system (47) obtained from $\mathcal{C}^{0}$ and $\mathcal{C}^{r}$ solvers for various diameters of initial conditions.


Figure 6: Plots of $t \rightarrow r\left(\left[D^{2} \phi\left(x_{0}, t\right)\right]\right)$ and $t \rightarrow r\left(\left[D^{3} \phi\left(x_{0}, t\right)\right]\right)$ for the Rössler system (49) obtained from $\mathcal{C}^{0}$ and $\mathcal{C}^{r}$ solvers for various diameters of initial conditions.


Figure 7: Plots of $t \rightarrow r\left(\left[D^{2} \phi\left(x_{0}, t\right)\right]\right)$ and $t \rightarrow r\left(\left[D^{3} \phi\left(x_{0}, t\right)\right]\right)$ for the HénonHeiles system (50) obtained from $\mathcal{C}^{0}$ and $\mathcal{C}^{r}$ solvers for various diameters of initial conditions.

| The system | second order derivatives |  |  | third order derivatives |  |  |
| :---: | ---: | ---: | ---: | ---: | ---: | ---: |
|  | $\mathcal{C}^{2}$-solver | $\mathcal{C}^{0}$-solver | ratio | $\mathcal{C}^{3}$-solver | $\mathcal{C}^{0}$-solver | ratio |
| V-L | 0.30 | 4.89 | 16 | 0.61 | 25.35 | 41 |
| pendulum | 0.09 | 1.96 | 21 | 0.21 | 9.71 | 46 |
| Michelson | 0.20 | 25.30 | 126 | 0.51 | 237.14 | 464 |
| Lorenz | 1.22 | 81.59 | 66 | 3.48 | 762.34 | 219 |
| Rössler | 0.71 | 58.96 | 83 | 1.87 | 521.37 | 278 |
| H-H | 1.40 | 430.21 | - | 4.96 | 3001.63 | - |

Table 1: Comparison of the time of computation of $\mathcal{C}^{2}$ and $\mathcal{C}^{3}$ solvers and of $\mathcal{C}^{0}$ solver when applied to the equations for variations. All the times of computation are given in seconds.

Comparison of the time of computations. As it is expected the $\mathcal{C}^{2}$ and $\mathcal{C}^{3}$ solvers are much faster than $\mathcal{C}^{0}$ applied to the equations for variations. In Table 1 we present the time of computation (in seconds) for each problem when computed from a point initial condition $\left(\operatorname{diam}\left(x_{0}\right)=0\right)$. For 2-3 dimensional systems the speed up of the computation of second order derivatives was between 16 and 126. For the third order derivatives it is even larger and varies between 41 and 464.

For the Hénon-Heiles Hamiltonian the $\mathcal{C}^{0}$-solver was not able to integrate along the periodic solution neither second nor third order derivatives even when starting from a point initial condition. In Table 1 we put gathered the times of computation up to the blow-up which occurred at $t=8.32874$ for the second order derivatives and $t=3.6712$ for the third order derivatives. The total time of integration for this system has been set to $T=13$.

Comparison of the obtained enclosures. For an interval $x=[a, b]$ we define a function

$$
r(x)=-\log _{10}\left(\frac{b-a}{|\operatorname{mid}(x)|}\right)=-\log _{10}\left(\frac{2(b-a)}{|a+b|}\right) .
$$

For an interval $x=[a, b]$ that does not contain zero, the function $r$ measures a relative diameter of $x$, i.e. an approximate number of significant decimal digits that are the same for $a$ and $b$.

With some abuse of notation we will denote by the same letter a relative diameter of an interval vector $[u] \subset \mathbb{R}^{m}$, i.e.

$$
r([u])=\min \left\{r\left([u]_{i}\right): i=1, \ldots, m\right\}
$$

and of an enclosure of $k$-th order derivative of a smooth function $f$

$$
r\left(\left[D^{k} f(x)\right]\right)=\min \left\{r\left(\left[D_{a} f(x)\right]\right):|a|=k\right\} .
$$

In Figures $2-7$ we present plots of the relative diameters of $r\left(D^{2} \phi\left(x_{0}, t\right)\right)$ and $r\left(D^{3} \phi\left(x_{0}, t\right)\right)$ as a function of time $t$ obtained from $\mathcal{C}^{0}$ and $\mathcal{C}^{r}$ solvers for various diameters of initial conditions and for the systems (45-50). Here $\phi$ denotes the local flow induced by the equation under consideration.

In principle, our $\mathcal{C}^{r}$-algorithm may be less accurate than the $\mathcal{C}^{0}$-Lohner direct solver in the computation of $D_{a} \varphi$ for $|a| \geq 1$, because we do not make use of the dependence of $D_{a} \varphi$ on $x$. Indeed, this can be seen for lower dimensional systems. But we have paid for this with the serious increase of the computation times.

For point initial conditions this lack of accuracy can be compensated by switching to the multiprecision arithmetic. In fact, for the systems under consideration we were able to obtain much thinner enclosures for derivatives using higher precision and within comparable or better time of computations.

For the initial conditions of nonzero diameters one can subdivide the sets. In many cases this strategy allows us to obtain better accuracy within same or better time. In some cases, like the Volterra-Lotka system (45) and the Lorenz system (47) obtained enclosures are significantly better when integrating the variational equations using the $\mathcal{C}^{0}$ solver. For these systems and low order derivatives one can choose between $\mathcal{C}^{0}$ and $\mathcal{C}^{r}$ solvers depending on the required accuracy of the result.

On the other hand, the $\mathcal{C}^{0}$ solver were not able to integrate third order derivatives for the Lorenz and Michelson systems when $\operatorname{diam}\left(x_{0}\right)=10^{-6}$.

For higher dimensional systems, like the Hénon-Heiles Hamiltonian we see that the $\mathcal{C}^{0}$-solver cannot compete with $\mathcal{C}^{r}$-algorithm. Both, the time of computation and obtained enclosures for second and third order derivatives are worse than those resulting from the $\mathcal{C}^{r}$-solver.

Memory usage. We would like to mention that the direct $\mathcal{C}^{0}$-solver when applied to the equations for variations requires also a huge memory. This is due to the fact, that the $\mathcal{C}^{0}$ solver extended by $k$-th order variational equations builds a tree for automatic differentiation for the system in $n\binom{n+k}{k}$-dimensional space and also its derivative. This squares the effective dimension for $\mathcal{C}^{0}$ solver. For the Hénon-Heiles system (50) and the third order derivatives the $\mathcal{C}^{0}$ solver used 22.31 MB of RAM while $\mathcal{C}^{3}$ solver used 416 kB , only.

Conclusions. The proposed algorithm has been proved to be very useful in many applications [KWZ, Wi, WZ2, WZ3]. In these papers we applied our $\mathcal{C}^{2}-\mathcal{C}^{5}$-algorithms to study various kind of dynamic and bifurcations of ODEs. In all of these applications the desired accuracy of computed derivatives was not that large as we usually require for the $\mathcal{C}^{0}$ image of the set - only a very few significant digits were necessary to get the result. Our tests show that $\mathcal{C}^{r}$ solver can compute high order derivatives with acceptable accuracy in a very good CPU time.

Our tests show also, that when the high accuracy of derivatives is required,
the $\mathcal{C}^{0}$ solver applied to the equations for variations can compete with $\mathcal{C}^{r}$ solvers for low dimensional systems and for low order derivatives, only. This is due to:

- loss of control of the wrapping effect in the $\mathcal{C}^{0}$ solver when the dimension is really high,
- memory usage. For example, using our $\mathcal{C}^{5}$ solver we integrated along a periodic solution the fifth order derivatives of a Hamiltonian flow (n-body problem) in 8 dimensions. The program used 7 GB of RAM. We were not able to build the $\mathcal{C}^{0}$ solver for fifth order derivatives on a computer with 64 GB of memory.
- Even if possible to build the necessary objects in the memory, the time of computations for large problems would be very large.


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