# Approximate invariant manifolds up to exponentially small terms 

Gérard Iooss ${ }^{1} \quad$ Eric Lombardi ${ }^{2}$<br>${ }^{1}$ I.U.F., Université de Nice, Labo J.A.Dieudonné, Parc Valrose, 06108 Nice, France<br>${ }^{2}$ Université Paul Sabatier, Institut de Mathématiques, 31062 Toulouse, France<br>gerard.iooss@unice.fr, lombardi@math.univ-toulouse.fr


#### Abstract

This paper is devoted to analytic vector fields near an equilibrium for which the linearized system is split in two invariant subspaces $E_{0}$ $\left(\operatorname{dim} m_{0}\right), E_{1}\left(\operatorname{dim} m_{1}\right)$. Under light diophantine conditions on the linear part, we prove that there is a polynomial change of coordinate in $E_{1}$ allowing to eliminate, in the $E_{1}$ component of the vector field, all terms depending only on the coordinate $u_{0} \in E_{0}$, up to an exponentially small remainder. This main result enables to prove the existence of analytic center manifolds up to exponentially small terms and extends to infinite dimensional vector fields. In the elliptic case, our results also proves, with very light assumptions on the linear part in $E_{1}$, that for initial data very close to a certain analytic manifold, the solution stays very close to this manifold for a very long time, which means that the modes in $E_{1}$ stay very small.

Keywords: analytic vector fields; normal forms; exponentially small remainders; center manifolds


AMS: 34M45; 34G20

## 1 Introduction

Let us consider an analytic vector field in the neighborhood of an equilibrium which we take at the origin. A natural idea is to try to uncouple a subset of coordinates from the other ones, by using a change of variables. This is used in particular since Poincaré and Dulac, and this is one of the main tool in the search of invariant manifolds of vector fields. Eliminating most of components of the vector field, expecting to only keep the relevant ones for the dynamics, is precisely the idea of center manifold reduction, which
is widely used in many physical systems, to simplify the study of the dynamics. However this reduction is only valid when we want to eliminate the hyperbolic part of the vector field and it has the defect to kill the analyticity after the reduction process. For systems fully elliptic near the origin, it may be expected to use a change of variables to uncouple all oscillatory modes. If this were possible, and if the initial data does not excite some modes, these ones would not be awaken for all times. Unfortunately, this is not possible in general, even though for hamiltonian systems, with suitable non resonant eigenvalues of the linearized system, it is nearly the case (Arnold diffusion between invariant tori corresponding to the "normal form" system with uncoupled modes).

In the present work, we consider systems for which the linearized system is split in two invariant subspaces $E_{0}\left(\operatorname{dim} m_{0}\right), E_{1}\left(\operatorname{dim} m_{1}\right)$. With light assumptions on the linear part, our main result is that there is a polynomial change of coordinate in $E_{1}$ allowing to eliminate, in the $E_{1}$ component of the vector field, all terms depending only on the coordinate $u_{0} \in E_{0}$, up to an exponentially small remainder (see Theorem 1). The proof of this theorem is based on a Gevrey estimate of the divergence of the remainder, which can be exponentially small by an optimal choice of the degree of the polynomial change of coordinates.

Gevrey estimates of the divergence of remainders, to get exponentially small upper bounds after an optimal choice of the order, were already used in the theory of normal forms for Hamiltonian systems in action-angle coordinates [2], [3], [14] following the pioneering work of Nekhoroshev [11, 12]. A similar result of exponential smallness of the remainder was also obtained by Giorgilli and Posilicano in [4] for a reversible system with a linear part composed of harmonic oscillators. For an extension of the result of normal forms with an exponentially remainder to any analytic vector fields with semi simple linear part see [7].

Direct normalization up to exponentially small terms is not available for vector fields studied in this paper since $\mathbf{L}_{1}$ is not assumed to be diagonalizable. However we can eliminate from the $E_{1}$ component of the vector field all terms depending only on the coordinate $u_{0} \in E_{0}$, up to an exponentially small remainder.

A first application of this result is when the linear part in $E_{1}$ is hyperbolic, while the linear part in $E_{0}$ has all its eigenvalues on the imaginary axis. It is well known that the center manifold reduction applies for small bounded solutions [8], which then lie on a manifold of same dimension as $E_{0}$. It is also well known that this manifold is in general not analytic [13], [20], [1], [16]. Our result allows to obtain a center manifold which is the graph of
a function sum of a polynomial of degree $p=O\left(\delta^{-b}\right)$ and an exponentially small function of order $O\left(e^{-c / \delta^{b}}\right)$ where $\delta$ is the size of the ball where we study the solutions, and $c$ and $b$ are positive numbers (see Theorem 5). It results in particular that the loss of analyticity is located in exponentially small terms. This result extends in infinite dimensional cases, then applicable in particular for a large class of PDE's. So combining, this result on center manifolds with the normal form theorem with exponentially small remainder [7] for the $E_{0}$ component ( $\mathbf{L}_{0}$ is diagonalizable), we can transform (1) into a new system with a "simplified" analytic leading part, perturbed by exponentially small terms. Such a transformation can be very useful when dealing with exponentially small phenomena (see [10]).

Another application, important in particular for engineering systems, is when the two linear subsystems in $E_{0}$ and $E_{1}$ have their eigenvalues on the imaginary axis. In particular, this situation happens for non linear vibrations of structures. Our result gives a sort of justification of a popular elimination process made in a formal way (see for example [9], [15], [17]), which allows to roughly state that for a class of initial data which do not excite in some sense the high frequencies (corresponding to $E_{1}$ ), then these ones are not awaken for all times....Our results prove, with very light diophantine assumptions (4) on the linear part in $E_{1}$, that for initial data very close to a certain analytic manifold, the solution stays very close to this manifold for a very long time, which means that the modes in $E_{1}$ stay very small (see theorem 8). This type of result is related to Arnold's diffusion for Hamiltonian systems (see a related result in [5]), while it should be noticed that we do not assume our system to be Hamiltonian, our assumptions on the eigenvalues being much lighter that usually done on such systems. In particular the linear part in $E_{1}$ is not assumed to be diagonalizable. Finally, notice the particular case studied in the same spirit by Groves and Schneider [6], for which $E_{0}$ is 2-dimensional and corresponds to a double eigenvalue in 0 , while $E_{1}$ corresponds to eigenvalues all imaginary. In this example there is no need of the diophantine condition (4), but $E_{1}$ is infinite dimensional and the result we obtain here needs to be adapted. In [18], Touzé and Amabili consider the damped case with an external periodic forcing. They assume that high frequency modes lie at a growing distance from the imaginary axis. Our method might be used in such a case, to rigorously prove that the high frequency modes do not awake as $t$ goes to infinity, provided certain non resonance condition between the forcing frequency and natural frequencies are realized, and provided the initial data is taken on a certain manifold in the spirit of Theorem 8 .

## 2 Main results

We gather in this section the main theorems proved in this paper. Our main theorem is the following

SplittingThm Theorem 1 Consider the following system in $\mathbb{R}^{m}$ (resp. $\mathbb{C}^{m}$ )

$$
\begin{equation*}
\frac{d u}{d t}=\mathbf{L} u+\mathbf{R}(u), \tag{1}
\end{equation*}
$$

where $u(t) \in \mathbb{R}^{m}$ (resp. $\left.\mathbb{C}^{m}\right), \mathbf{L}$ is a linear operator, and $\mathbf{R}$ is analytic in a neighborhood of the origin, such that

$$
\begin{equation*}
\mathbf{R}(u)=\sum_{2 \leq k} \mathbf{R}_{k}\left[u^{(k)}\right] \tag{2}
\end{equation*}
$$

where $\mathbf{R}_{k}$ is a $k$ - linear symmetric map on $\left(\mathbb{R}^{m}\right)^{k}\left(\right.$ resp. $\left.\left(\mathbb{C}^{m}\right)^{k}\right)$ satisfying

$$
\begin{equation*}
\left\|\mathbf{R}_{k}\left[u_{1}, u_{2}, \cdots, u_{k}\right]\right\| \leq \frac{c}{\rho^{k}}\left\|u_{1}\right\| \cdots\left\|u_{k}\right\|, \tag{3}
\end{equation*}
$$

for a certain radius of convergence $\rho>0$ (here $\left[u^{(k)}\right]$ means the $k$-uple of vectors $[u, u, \cdots, u])$. Assume that the linear operator $\mathbf{L}$ is the direct sum of two linear operators $\mathbf{L}_{0}$ on $E_{0}\left(\operatorname{dim} m_{0}\right)$, and $\mathbf{L}_{1}$ on $E_{1}\left(\operatorname{dim} m_{1}\right)$, such that $\mathbf{L}_{0}$ is diagonalizable with eigenvalues $\lambda_{1}^{(0)}, \cdots, \lambda_{m_{0}}^{(0)}$ and that there exist constants $\gamma>0, \tau \geq 0$ such that

$$
\begin{equation*}
\left|\left\langle\alpha, \lambda^{(0)}\right\rangle-\lambda_{j}^{(1)}\right| \geq \frac{\gamma}{|\alpha|^{\tau}} \tag{4}
\end{equation*}
$$

holds for any $\alpha \in \mathbb{N}^{m_{0}} \backslash\{0\}$, and any eigenvalue $\lambda_{j}^{(1)}$ of $\mathbf{L}_{1}$.
Then there exists a polynomial $\boldsymbol{\Phi}: E_{0} \rightarrow E_{1}$ of degree $p=O\left(\delta^{-b}\right)$ such that the change of variables in $E_{1}$

$$
\begin{equation*}
u_{1}=v_{1}+\boldsymbol{\Phi}\left(u_{0}\right) \tag{5}
\end{equation*}
$$

Change Var
transforms the system (1) into the following system in $E_{0} \times E_{1}$ :

$$
\begin{align*}
\frac{d u_{0}}{d t} & =\mathbf{L}_{0} u_{0}+\mathbf{R}^{(0)}\left(u_{0}, v_{1}\right),  \tag{6}\\
\frac{d v_{1}}{d t} & =\mathbf{L}_{1} v_{1}+\mathbf{R}^{(1)}\left(u_{0}, v_{1}\right)+\boldsymbol{\rho}\left(u_{0}\right),
\end{align*}
$$

in which $\mathbf{R}^{(0)}, \mathbf{R}^{(1)}, \boldsymbol{\rho}$ are analytic in their arguments, and where

$$
\mathbf{R}^{(0)}\left(u_{0}, u_{1}\right)=\mathbf{P}_{0} \mathbf{R}\left(u_{0}+v_{1}+\boldsymbol{\Phi}\left(u_{0}\right)\right),
$$

$\mathbf{P}_{0}$ being the projection on $E_{0}$ which commutes with $\mathbf{L}$, and

$$
\begin{gather*}
\mathbf{R}^{(1)}\left(u_{0}, v_{1}\right)=O\left(\left\|v_{1}\right\|\left(\left\|u_{0}\right\|+\left\|v_{1}\right\|\right)\right),  \tag{7}\\
\sup _{\left\|u_{0}\right\| \leq \delta}\left\|\boldsymbol{\rho}\left(u_{0}\right)\right\| \leq M e^{-\frac{w}{\delta^{b}}}, \tag{8}
\end{gather*}
$$

with $M, w>0$ depending only on $\tau, m_{0}, c, \rho, \mathbf{L}_{1}$ and

$$
b=\frac{1}{1+\nu \tau}
$$

where $\nu$ is the maximal index (size of Jordan blocks) of eigenvalues of $\mathbf{L}_{1}$.
Remark 2 Notice that the constants $M$ and $w$ do not depend on the dimension $m_{1}$ of the subspace $E_{1}$ if $\mathbf{L}_{1}$ is a priori in Jordan form. This allows to consider systems with large (even infinite) dimensions.

Remark 3 Since all the norms are equivalent on $\mathbb{R}^{m}$ (resp. $\mathbb{C}^{m}$ ), (7),(8) remains true for any norm on $\mathbb{R}^{m}$ (resp. $\mathbb{C}^{m}$ ). A change of norm simply change the values of $M$ and $w$. So, estimates (7),(8) remain true under linear change of coordinates up to a change of values of $M$ and $w$. Hence without loss of generality we can assume that the complexified space of $E_{0}$ and $E_{1}$, still denoted by $E_{0}$ and $E_{1}$ read respectively $E_{0}=\mathbb{C}^{m_{0}} \times \underbrace{(0, \cdots, 0)}_{m_{1} \text { times }}$ and $E_{1}=\underbrace{(0, \cdots, 0)}_{m_{0} \text { times }} \times \mathbb{C}^{m_{1}}$ and that in the canonical basis of $\mathbb{C}^{m}, L_{0}$ is diagonal and $L_{1}$ is under Jordan normal form.

We deduce from the above theorem a corollary which deals with vector fields depending on parameters.

Corollary 4 Consider the following system in $\mathbb{R}^{m}$ (resp. $\mathbb{C}^{m}$ )

$$
\begin{equation*}
\frac{d u}{d t}=\mathbf{L} u+\mathbf{R}(u, \mu), \tag{9}
\end{equation*}
$$

where $u(t) \in \mathbb{R}^{m}$ (resp. $\mathbb{C}^{m}$ ), $\mathbf{L}$ is a linear operator, and $\mathbf{R}$ is analytic in a neighborhood of the origin in $\mathbb{R}^{m} \times \mathbb{R}^{q}$ (resp. $\mathbb{C}^{m} \times \mathbb{R}^{q}$ ) and such that

$$
\begin{equation*}
\mathbf{R}(0, \mu)=0, D_{u} \mathbf{R}(0,0)=0 . \tag{10}
\end{equation*}
$$

Assuming the same hypothesis on $\mathbf{L}$ as in Theorem 1, and that 0 is not eigenvalue of $\mathbf{L}_{1}$, then there exists a polynomial $\boldsymbol{\Phi}: E_{0} \times \mathbb{R}^{q} \rightarrow E_{1}$ of degree $p=O\left(\delta^{-b}\right)$ such that the change of variables in $E_{1}$

$$
u_{1}=v_{1}+\boldsymbol{\Phi}\left(u_{0}, \mu\right)
$$

transforms the system (1) into the following system in $E_{0} \times E_{1}$ :

$$
\begin{aligned}
\frac{d u_{0}}{d t} & =\mathbf{L}_{0} u_{0}+\mathbf{R}^{(0)}\left(u_{0}, v_{1}, \mu\right) \\
\frac{d v_{1}}{d t} & =\mathbf{L}_{1} v_{1}+\mathbf{R}^{(1)}\left(u_{0}, v_{1}, \mu\right)+\boldsymbol{\rho}\left(u_{0}, \mu\right)
\end{aligned}
$$

in which $\mathbf{R}^{(0)}, \mathbf{R}^{(1)}, \boldsymbol{\rho}$ are analytic in their arguments, and where

$$
\mathbf{R}^{(0)}\left(u_{0}, u_{1}, \mu\right)=\mathbf{P}_{0} \mathbf{R}\left(u_{0}+v_{1}+\boldsymbol{\Phi}\left(u_{0}, \mu\right), \mu\right)
$$

$\mathbf{P}_{0}$ being the projection on $E_{0}$ which commutes with $\mathbf{L}$, and

$$
\begin{gathered}
\mathbf{R}^{(1)}\left(u_{0}, v_{1}, \mu\right)=O\left(\left\|v_{1}\right\|\left(\left\|u_{0}\right\|+\left\|v_{1}\right\|+\|\mu\|\right)\right) \\
\sup _{\left\|u_{0}\right\|+\|\mu\| \leq \delta}\left\|\boldsymbol{\rho}\left(u_{0}, \mu\right)\right\| \leq M e^{-\frac{w}{\delta^{b}}}
\end{gathered}
$$

with $M, w>0$ depending only on $\tau, m_{0}, c, \rho, \mathbf{L}_{1}$ and $b$ is as in Theorem 1.
Another application of theorem 1, is the existence of analytic center manifolds up to exponentially small term. More precisely, consider the case when the spectrum of $\mathbf{L}_{0} \subset i \mathbb{R}$, and $\mathbf{L}_{1}$ is hyperbolic, i.e. the eigenvalues of $\mathbf{L}_{1}$ lie at a distance $\gamma>0$ from the imaginary axis. Then in finite dimension we have the following

## Theorem 5 (Center manifold analytic up to exp. small terms)

Consider the analytic system (1) in $\mathbb{R}^{m}$ and assume that $\mathbf{L}_{0}$ is diagonalizable with all its eigenvalues on the imaginary axis, and assume that $\mathbf{L}_{1}$ has its eigenvalues at least at a distance $\gamma>0$ from the imaginary axis.

Then for any $k \geq 2$, there exists a polynomial $\mathbf{\Phi}: E_{0} \rightarrow E_{1}$ of degree $O(1 / \delta)$, with $\boldsymbol{\Phi}(0)=0, D \boldsymbol{\Phi}(0)=0$, a neighborhood $\mathcal{O}$ of 0 in $\mathbb{R}^{m}$, and a $\operatorname{map} \boldsymbol{\Psi} \in \mathcal{C}^{k}\left(E_{0}, E_{1}\right)$ which is $O\left(e^{-\frac{C}{\delta}}\right)$ for $\left\|u_{0}\right\|_{E_{0}} \leq \delta$ and a certain constant $C>0$, such that the manifold

$$
\begin{equation*}
\mathcal{M}_{0}=\left\{u_{0}+\boldsymbol{\Phi}\left(u_{0}\right)+\boldsymbol{\Psi}\left(u_{0}\right) ; u_{0} \in E_{0}\right\} \tag{11}
\end{equation*}
$$

has the following properties.
(a) $\mathcal{M}_{0}$ is locally invariant, i.e., if $u$ is a solution of (1) satisfying $u(0) \in$ $\mathcal{M}_{0} \cap \mathcal{O}$ and $u(t) \in \mathcal{O}$ for all $t \in[0, T]$, then $u(t) \in \mathcal{M}_{0}$ for all $t \in[0, T]$.
(b) $\mathcal{M}_{0}$ contains the set of bounded solutions of (1) staying in $\mathcal{O}$ for all $t \in \mathbb{R}$, i.e., if $u$ is a solution of (1) satisfying $u(t) \in \mathcal{O}$ for all $t \in \mathbb{R}$, then $u(0) \in \mathcal{M}_{0}$.

Remark 6 The interest of Theorem 5 is that it implies that the reduced system on the center manifold is analytic, up to exponentially small terms. This property is clearly still true after the polynomial new change of variables which put the reduced system under normal form (the usual one). In considering the analytic part of the reduced vector field, this normal form may be derived up to an optimal degree, as made in [7], since $\mathbf{L}_{0}$ is diagonalizable. This may be helpful when dealing with exponentially small phenomena associated with the original system (1).

Remark 7 This theorem is also true in the infinite dimensional case (see Theorem 20 in subsection 4.2)

A last application of theorem 1, important in particular for engineering systems, is when the two linear subsystems in $E_{0}$ and $E_{1}$ have both their eigenvalues on the imaginary axis. More precisely in section 5 , we prove

EllipticThm Theorem 8 (Elliptic vector fields) Assume that assumptions of Theorem 1 hold, and in addition that $\mathbf{L}_{1}$ has only imaginary eigenvalues. Then for any small initial data $u(0)$ chosen on the manifold $\mathcal{M}_{0}^{\prime}=\left\{u=u_{0}+\boldsymbol{\Phi}\left(u_{0}\right) ; u_{0} \in\right.$ $\left.E_{0}\right\}$ the solution $u(t)$ stays at a distance $O\left(e^{-\frac{C}{\delta^{b}}}\right)$ to $\mathcal{M}_{0}^{\prime}$ for $t \in[0, T]$, with $T=O\left(\delta^{-(b+1 / \nu)}\right)$, where $b=(1+\nu \tau)^{-1}$ and $\nu$ is the maximal index of eigenvalues of $\mathbf{L}_{1}$.

Remark 9 We observe (see 39) that in going up to exponentially small terms in Theorem 1, we win the exponential smallness of $\left\|v_{1}(t)\right\|$ for a long range of time, without more precise assumption on $\mathbf{L}_{1}$. If we assume more specific properties of the system, we may have a longer range of time for the validity of this exponential smallness. First, if $\mathbf{L}_{1}$ is diagonalizable this range of time is $O\left(\delta^{\left[1+(1+\tau)^{-1}\right]}\right)$.

Remark 10 Let assume in addition that (1) is a reversible system such that $\mathbf{L}$ has only pairs of simple imaginary eigenvalues, satisfying the $\gamma, \tau$ homologically diophantine assumption defined in [7]: for every $\alpha \in \mathbb{N}^{m},|\alpha| \geq$ 2

$$
\left|\langle\alpha, \lambda\rangle-\lambda_{j}\right| \geq \frac{\gamma}{|\alpha|^{\tau}} \text { when }\langle\alpha, \lambda\rangle-\lambda_{j} \neq 0
$$

and $\langle\alpha, \lambda\rangle-\lambda_{j}=0$ only for the trivial cases $2 \lambda_{j}+\overline{\lambda_{j}}-\lambda_{j}=0$ (non resonance assumption). In such a case, we can use the normal form theorem of [7] which gives a normal form up to an exponentially small term, which improves the final form of Theorem 1 since the coupling between the subsystems in $E_{0}$ and in $E_{1}$ only appears in exponentially small terms. Taking $v_{1}(0)=0$, it is easy to show that $\left\|v_{1}(t)\right\|$ stays exponentially small now for an exponentially long time (analogue to Arnold diffusion).

## 3 Proof of the main theorem

We first deduce corollary 4 from theorem 1 and then we prove this theorem.
Proof of Corollary 4. Let us define

$$
\widetilde{u}=(u, \mu) \in \mathbb{R}^{m} \times \mathbb{R}^{q}
$$

then the system reads

$$
\begin{equation*}
\frac{d \widetilde{u}}{d t}=\widetilde{\mathbf{L}} \widetilde{u}+\widetilde{\mathbf{R}}(\widetilde{u}), \tag{12}
\end{equation*}
$$

with

$$
\widetilde{\mathbf{L}} \widetilde{u}=(\mathbf{L} u, 0), \widetilde{\mathbf{R}}(\widetilde{u})=(\mathbf{R}(u, \mu), 0)
$$

Then, it is clear that the system (12) satisfies all assumptions of Theorem 1. In particular, the operator $\widetilde{\mathbf{L}}$ is the direct sum of $\widetilde{\mathbf{L}}_{0}$ and $\widetilde{\mathbf{L}}_{1}$ defined by

$$
\begin{aligned}
\widetilde{\mathbf{L}}_{0} \widetilde{u}_{0} & =\left(\mathbf{L}_{0} u_{0}, 0\right), \text { for } \widetilde{u}_{0} \in \widetilde{E_{0}}=E_{0} \times \mathbb{R}^{q} \\
\widetilde{\mathbf{L}}_{1} \widetilde{u}_{1} & =\left(\mathbf{L}_{1} u_{1}, 0\right), \text { for } \widetilde{u}_{1} \in \widetilde{E_{1}}=E_{1} \times\{0\},
\end{aligned}
$$

and the eigenvalues of $\widetilde{\mathbf{L}}_{1}$ are those of $\mathbf{L}_{1}$, while the eigenvalues of $\widetilde{\mathbf{L}}_{0}$ are those of $\mathbf{L}_{0}$ with 0 still semi-simple, having an additional $q$-dimensional eigenspace: $(0, \mu), \mu \in \mathbb{R}^{q}$ and the diophantine condition (4) is still satisfied. Hence the Corollary is proved.

Proof of Theorem 1. In the proof below we use several algebraic properties which were proved in [7]. Performing the change of coordinates $u=u_{0}+u_{1}+\Phi\left(u_{0}\right)$, we check that (1) is equivalent to (6) close to the origin if and only if

$$
\begin{aligned}
\boldsymbol{P}_{\mathbf{0}} \boldsymbol{R}\left(u_{0}+v_{1}+\boldsymbol{\phi}\left(u_{0}\right)\right)= & \boldsymbol{R}^{(0)}\left(u_{0}, v_{1}\right) \\
D \boldsymbol{\Phi}\left(u_{0}\right) \cdot \boldsymbol{L}_{0} u_{0}-\boldsymbol{L}_{1} \boldsymbol{\Phi}\left(u_{0}\right)= & -D \boldsymbol{\Phi}\left(u_{0}\right) \cdot \boldsymbol{R}^{(0)}\left(u_{0}, v_{1}\right)-\boldsymbol{\rho}\left(u_{0}\right) \\
& +\boldsymbol{P}_{\mathbf{1}} \boldsymbol{R}\left(u_{0}+u_{1}+\boldsymbol{\Phi}\left(u_{0}\right)\right)-\boldsymbol{R}^{(1)}\left(u_{0}, v_{1}\right) .
\end{aligned}
$$

Then Setting $v_{1}=0$ and using (7), we obtain the following basic identity

$$
\begin{align*}
& D \boldsymbol{\Phi}\left(u_{0}\right) \mathbf{L}_{0} u_{0}-\mathbf{L}_{1} \boldsymbol{\Phi}\left(u_{0}\right)=-D \boldsymbol{\Phi}\left(u_{0}\right) \mathbf{P}_{0} \\
& \mathbf{R}\left(u_{0}+\boldsymbol{\Phi}\left(u_{0}\right)\right)  \tag{13}\\
&+\mathbf{P}_{1} \mathbf{R}\left(u_{0}+\boldsymbol{\Phi}\left(u_{0}\right)\right)-\boldsymbol{\rho}\left(u_{0}\right) .
\end{align*}
$$

Let decompose the polynomial $\boldsymbol{\Phi}$ into a sum of homogeneous polynomials of increasing degrees

$$
\boldsymbol{\Phi}\left(u_{0}\right)=\sum_{2 \leq k \leq p} \boldsymbol{\Phi}_{k}\left[u_{0}^{(k)}\right]
$$

with $k$ - linear symmetric maps $\boldsymbol{\Phi}_{k}:\left(E_{0}\right)^{k} \rightarrow E_{1}$. For convenience we denote by $\boldsymbol{\Phi}_{1}\left(u_{0}\right) \equiv u_{0}$ which takes its values in $E_{0}$ (contrary to $\boldsymbol{\Phi}_{k}$ for $k \geq 2$, which takes its values in $E_{1}$ ). Then we have for $2 \leq n \leq p$

$$
\begin{equation*}
D \boldsymbol{\Phi}_{n}\left[u_{0}^{(n)}\right] \mathbf{L}_{0} u_{0}-\mathbf{L}_{1} \boldsymbol{\Phi}_{n}\left[u_{0}^{(n)}\right]=\mathbf{F}_{n}\left[u_{0}^{(n)}\right] \tag{14}
\end{equation*}
$$

with

$$
\begin{aligned}
\mathbf{F}_{n}\left[u_{0}^{(n)}\right]= & \sum_{\substack{2 \leq q \leq n \\
k_{1}+\cdots+k_{q}=n, k_{j} \geq 1}} \mathbf{P}_{1} \mathbf{R}_{q}\left[\mathbf{\Phi}_{k_{1}}, \cdots, \boldsymbol{\Phi}_{k_{q}}\right]+ \\
& -\sum_{\substack{2 \leq \ell \leq n-1,2 \leq q \leq n-\ell+1 \\
k_{1}+\cdots+k_{q}=n-\ell+1, k_{j} \geq 1}} D \boldsymbol{\Phi}_{\ell}\left[u_{0}^{(\ell)}\right] \mathbf{P}_{0} \mathbf{R}_{q}\left[\boldsymbol{\Phi}_{k_{1}}, \cdots, \boldsymbol{\Phi}_{k_{q}}\right] .
\end{aligned}
$$

Equation (14) is of the form

$$
\mathcal{A} \mathbf{\Phi}_{n}=\mathbf{F}_{n}
$$

with the homological operator $\mathcal{A}$ defined on the vector space of polynomials $\boldsymbol{\Phi}: E_{0} \rightarrow E_{1}$, by

$$
\begin{equation*}
\mathcal{A} \boldsymbol{\Phi}=D \boldsymbol{\Phi}\left(u_{0}\right) \mathbf{L}_{0} u_{0}-\mathbf{L}_{1} \boldsymbol{\Phi}\left(u_{0}\right) \tag{15}
\end{equation*}
$$

We then need to introduce the scalar product in the space $\mathcal{H}$ of polynomials of a variable in $E_{0}$, taking values in $\mathbb{C}^{m}$ (which could be in the complexified space of the subspace $E_{1}$ or $E_{0}$ ) as done in [7].

Given two polynomials $\boldsymbol{\Phi}$ and $\boldsymbol{\Phi}^{\prime}$ we define their scalar product by

$$
\left\langle\boldsymbol{\Phi}, \boldsymbol{\Phi}^{\prime}\right\rangle_{\mathcal{H}}:=\sum_{1 \leq j \leq n}\left\langle\boldsymbol{\Phi}_{j}, \boldsymbol{\Phi}_{j}^{\prime}\right\rangle
$$

with $\boldsymbol{\Phi}=\left(\mathbf{\Phi}_{1}, \cdots, \mathbf{\Phi}_{n}\right), \boldsymbol{\Phi}^{\prime}=\left(\mathbf{\Phi}_{1}^{\prime}, \cdots, \boldsymbol{\Phi}_{n}^{\prime}\right)$, and where for a pair of polynomials $P, Q: E_{0} \rightarrow \mathbb{C}$,

$$
\langle P, Q\rangle=\left.\bar{P}\left(\partial_{X}\right) Q(X)\right|_{X=0}
$$

where by definition

$$
\bar{P}(X)=\overline{P(\bar{X})}
$$

Then the associated euclidian norm is defined by

$$
|\boldsymbol{\Phi}|_{2}:=\sqrt{\langle\boldsymbol{\Phi}, \boldsymbol{\Phi}\rangle_{\mathcal{H}}} .
$$

It is clear that for any $n \geq 1$, the linear operator $\mathcal{A}$ leaves invariant the subspace $\mathcal{H}_{n}$ of homogeneous polynomials of degree $n$, and we have the following Lemma proved in Appendix:

Eric's Lemma Lemma 11 The operator $\mathcal{A}$ is invertible in the subspace $\mathcal{H}_{n}$ and there exists a constant $a$, depending only on $\gamma$ and $\mathbf{L}_{1}$, such that

$$
\left|\left\|\left.|\mathcal{A}|_{\mathcal{H}_{n}}^{-1}\left|\|_{2}:=\sup _{|\boldsymbol{\Phi}|_{2}=1}\right| \mathcal{A}\right|_{\mathcal{H}_{n}} ^{-1} \boldsymbol{\Phi} \mid \leq a n^{\tau^{\prime}}\right.\right.
$$

where $\tau^{\prime}=\nu \tau$, and $\nu$ is the maximal index of the eigenvalues of $\mathbf{L}_{1}$.
This lemma is proved in Appendix A.
Remark 12 If $\mathbf{L}_{1}$ is in Jordan form, the constant a depends only on $\gamma$ and $\nu$. If $\mathbf{L}_{1}$ is diagonal then $\tau^{\prime}=\tau$ and $a=1 / \gamma$.

Moreover, defining the norm

$$
\phi_{n}:=|\boldsymbol{\Phi}|_{2, n}:=\frac{1}{\sqrt{n!}}|\boldsymbol{\Phi}|_{2}, \text { for } \boldsymbol{\Phi} \in \mathcal{H}_{n}
$$

we have the following lemma, proved in [7] (see Lemmas 2.10, 2.11):

## Lemma 13

(i) For $k_{1}+\ldots+k_{q}=n$

$$
\left|\mathbf{R}_{q}\left[\mathbf{\Phi}_{k_{1}}, \cdots, \mathbf{\Phi}_{k_{q}}\right]\right|_{2, n} \leq \frac{c}{\rho^{q}} \phi_{k_{1}} \cdots \phi_{k_{q}}
$$

(ii) for $2 \leq \ell \leq p$, $\ell+k=n+1$, and any $\mathbf{N}_{k} \in \mathcal{H}_{k}$

$$
\left|D \mathbf{\Phi}_{\ell} \cdot \mathbf{N}_{k}\right|_{2, n} \leq \sqrt{\ell^{2}+\left(m_{0}-1\right) \ell} \phi_{\ell}\left|\mathbf{N}_{k}\right|_{2, k} \leq \ell \sqrt{m_{0}} \phi_{\ell}\left|\mathbf{N}_{k}\right|_{2, k}
$$

Then, the proof of Theorem 1 is performed in several steps giving respectively estimates of $\phi_{n},\left\|\sum \boldsymbol{\Phi}\left(u_{0}\right)\right\|$, and $\rho_{0}$ gathered in the following lemmas:

Lemphin Lemma 14 There exists $K>0$ depending only on $c, c_{01}, \rho, m_{0}$, a such that for every $n$ with $1 \leq n \leq p$,

$$
\begin{equation*}
\phi_{n} \leq \sqrt{m_{0}} K^{n-1}(n!)^{1+\tau^{\prime}} \tag{16}
\end{equation*}
$$

where $c_{01}:=\max \left(\left\|\left|\boldsymbol{P}_{0}\| \|,\left\|\mid \boldsymbol{P}_{1}\right\| \|\right)\right.\right.$.
Lemma 15 Let us choose $p$ such that

$$
\begin{equation*}
p=p_{\mathrm{opt}}:=\left[\frac{1}{(2 \delta K)^{b}}\right], \quad b=\frac{1}{1+\tau^{\prime}} \tag{17}
\end{equation*}
$$

where [•] denotes the integer part of a number. Then for $\left\|u_{0}\right\| \leq \delta$ we have

$$
\left\|\sum_{1 \leq k \leq p_{\mathrm{opt}}} \boldsymbol{\Phi}_{k}\left(u_{0}\right)\right\| \leq 2 \delta \sqrt{m_{0}}
$$

rhoopt Lemma 16 The remainder $\boldsymbol{\rho}$ satisfies

$$
\boldsymbol{\rho}\left(u_{0}\right)=\mathfrak{R}_{1}\left(u_{0}\right)+\mathfrak{R}_{2}\left(u_{0}\right)+\mathfrak{R}_{3}\left(u_{0}\right)+\mathfrak{R}_{4}\left(u_{0}\right),
$$

with

$$
\begin{aligned}
& \mathfrak{R}_{1}\left(u_{0}\right)=\sum_{p+1 \leq q} \mathbf{P}_{1} \mathbf{R}_{q}\left[\left(\sum_{1 \leq k \leq p} \boldsymbol{\Phi}_{k}\left(u_{0}\right)\right)^{(q)}\right], \\
& \mathfrak{R}_{2}\left(u_{0}\right)=-\sum_{2 \leq \ell \leq p, p+1 \leq q} D \boldsymbol{\Phi}_{\ell}\left[u_{0}^{(\ell)}\right] \mathbf{P}_{0} \mathbf{R}_{q}\left[\left(\sum_{1 \leq k \leq p} \boldsymbol{\Phi}_{k}\left(u_{0}\right)\right)^{(q)}\right], \\
& \mathfrak{R}_{3}\left(u_{0}\right)=\sum_{\substack{2 \leq q \leq p, 1 \leq k_{j} \leq p \\
k_{1}+\cdots+k_{q} \geq p+1}} \mathbf{P}_{1} \mathbf{R}_{q}\left[\boldsymbol{\Phi}_{k_{1}}\left(u_{0}\right), \cdots, \boldsymbol{\Phi}_{k_{q}}\left(u_{0}\right)\right], \\
& \mathfrak{R}_{4}\left(u_{0}\right)=-\sum_{\substack{2 \leq \ell \leq p, k_{1}+\ldots k_{q} \geq p-l+2}} D \boldsymbol{\Phi}_{\ell}\left[u_{0}^{(\ell)}\right] \mathbf{P}_{0} \mathbf{R}_{q}\left[\boldsymbol{\Phi}_{k_{1}}\left(u_{0}\right), \cdots, \boldsymbol{\Phi}_{k_{q}}\left(u_{0}\right)\right],
\end{aligned}
$$

and for $p=p_{\mathrm{opt}}$, it satisfies

$$
\begin{equation*}
\sup _{\left\|u_{0}\right\| \leq \delta}\left\|\boldsymbol{\rho}\left(u_{0}\right)\right\| \leq M e^{-\frac{w}{\delta^{b}}} \tag{18}
\end{equation*}
$$

with $M, w>0$ depending only on $\tau, m_{0}, c, \rho, \mathbf{L}_{1}$.

Proof of lemma 14. Estimate of $\phi_{n}$. We obtain from (14)

$$
\phi_{n} \leq a c c_{01} n^{\tau^{\prime}}\left\{\begin{array}{c}
\sum_{\substack{2 \leq q \leq n \\
k_{1}+\cdots+k_{q}=n, k_{j} \geq 1}} \frac{1}{\rho^{q}} \phi_{k_{1}} \cdots \phi_{k_{q}} \\
\left.+\sum_{\substack{2 \leq \ell \leq n-1,2 \leq q \leq n-\ell+1 \\
k_{1}+\cdots+k_{q}=n-\ell+1, k_{j} \geq 1}} \frac{\ell \sqrt{m_{0}}}{\rho^{q}} \phi_{\ell} \phi_{k_{1}} \cdots \phi_{k_{q}}\right\} . \tag{19}
\end{array}\right.
$$

Then, notice that by construction

$$
\phi_{1}=\sqrt{m_{0}}
$$

For suppressing the factor $n^{\tau^{\prime}}$ in the inequality (19), we introduce the following sequence $\alpha_{n}$ defined by

$$
\alpha_{1}=1 \quad \text { and } \quad \phi_{n}=\sqrt{m_{0}} K_{1}^{n-1}(n!)^{\tau^{\prime}} \alpha_{n}, \quad \text { for } \quad n \geq 1
$$

where $K_{1}$ will be chosen later. Using the following inequalities proved in [7]-lemma 2.12,

$$
\frac{k_{1}!\cdots+k_{q}!}{(n-1)!} \leq 1 \quad \text { for } 2 \leq q \leq n, k_{1}+\cdots+k_{q}=n
$$

and

$$
\begin{gathered}
\frac{\ell!k_{1}!\cdots k_{q}!}{(n-1)!}=\frac{\ell!k_{1}!\cdots k_{q}!}{(n-\ell)!} \frac{(n-\ell)!}{(n-1)!} \leq 1, \quad \text { for }\left\{\begin{array}{l}
2 \leq \ell \leq n-1 \\
2 \leq q \leq n-\ell+1 \\
k_{1}+\cdots+k_{q}=n-\ell+1,
\end{array}\right. \\
\alpha_{n} \leq \frac{a c c_{01} K_{1}}{\sqrt{m_{0}}}\left\{\begin{array}{l}
\sum_{\substack{2 \leq q \leq n \\
k_{1}+\cdots+k_{q}=n, k_{j} \geq 1}}\left(\frac{\sqrt{m_{0}}}{K_{1} \rho}\right)^{q} \alpha_{k_{1} \cdots \alpha_{k_{q}}+} \\
\sum_{\substack{2 \leq \ell \leq n-1,2 \leq q \leq n-\ell+1 \\
k_{1}+\cdots+k_{q}=n-\ell+1, k_{j} \geq 1}} m_{0}\left(\frac{\sqrt{m_{0}}}{K_{1} \rho}\right)^{q} \ell \alpha_{\ell} \alpha_{k_{1}} \cdots \alpha_{k_{q}}
\end{array}\right\}
\end{gathered}
$$

and by choosing

$$
\begin{equation*}
K_{1} \geq \frac{a c c_{01} m_{0}^{3 / 2}}{\rho^{2}} \tag{20}
\end{equation*}
$$

we finally get

$$
\begin{align*}
\alpha_{n} \leq & \sum_{\substack{2 \leq q \leq n \\
k_{1}+\cdots+k_{q}=n, k_{j} \geq 1}}\left(\frac{\sqrt{m_{0}}}{K_{1} \rho}\right)^{q-2} \alpha_{k_{1}} \cdots \alpha_{k_{q}}+  \tag{21}\\
& +\sum_{\substack{2 \leq \ell \leq n-1,2 \leq q \leq n-\ell+1 \\
k_{1}+\cdots+k_{q}=n-\ell+1, k_{j} \geq 1}}\left(\frac{\sqrt{m_{0}}}{K_{1} \rho}\right)^{q-2} \ell \alpha_{\ell} \alpha_{k_{1}} \cdots \alpha_{k_{q}}
\end{align*}
$$

Now, the idea is to use the majorizing sequence $\beta_{n}$ defined by

$$
\beta_{1}=1, \beta_{n}=\Theta^{n-2}(n-2)!\text { for } n \geq 2
$$

the number $\Theta$ being chosen later, large enough. It is clear that

$$
\begin{aligned}
& \alpha_{1}=1 \leq \beta_{1} \\
& \alpha_{2} \leq \alpha_{1}^{2}=1 \leq \beta_{2}
\end{aligned}
$$

Assuming that $\alpha_{k} \leq \beta_{k}$ for $1 \leq k \leq n-1$, we intend to prove that $\alpha_{n} \leq \beta_{n}$. Indeed, by replacing $\alpha_{k}$ by $\beta_{k}, 1 \leq k \leq n-1$, in the right hand side of (21) we find that

$$
\alpha_{n} \leq \sum_{2 \leq q \leq n}\left(\frac{\sqrt{m_{0}}}{K_{1} \rho}\right)^{q-2} \Pi_{q, n}+\sum_{2 \leq \ell \leq n-1,2 \leq q \leq n-\ell+1}\left(\frac{\sqrt{m_{0}}}{K_{1} \rho}\right)^{q-2} \ell \beta_{\ell} \Pi_{q, n-\ell+1}
$$

with, for $2 \leq q \leq n$

$$
\Pi_{q, n}:=\sum_{k_{1}+\cdots+k_{q}=n, k_{j} \geq 1} \beta_{k_{1}} \cdots \beta_{k_{q}}
$$

It is shown in [7]-lemma 2.13 that

$$
\begin{array}{ll}
\Pi_{2, n} \leq \frac{2}{\Theta} \beta_{n}, & \text { for } n \geq 3 \\
\Pi_{q, n} \leq \frac{2}{\Theta^{q-2}} \beta_{n}, & \text { for } 3 \leq q \leq n
\end{array}
$$

Hence

$$
\begin{aligned}
\alpha_{n} \leq & \left(\frac{2}{\Theta}+\sum_{3 \leq q \leq n} 2\left(\frac{\sqrt{m_{0}}}{\Theta K_{1} \rho}\right)^{q-2}\right) \beta_{n}+ \\
& +\frac{2}{\Theta} \sum_{2 \leq \ell \leq n-1} \ell \beta_{\ell} \beta_{n-\ell+1}+\sum_{2 \leq \ell \leq n-1,3 \leq q \leq n-\ell+1} 2\left(\frac{\sqrt{m_{0}}}{\Theta K_{1} \rho}\right)^{q-2} \ell \beta_{\ell} \beta_{n-\ell+1} .
\end{aligned}
$$

We choose now $\Theta$ and $K_{1}$ such that

$$
\begin{equation*}
\frac{1}{\Theta}+\frac{\sqrt{m_{0}}}{\Theta K_{1} \rho-\sqrt{m_{0}}} \leq \frac{1}{4} \tag{22}
\end{equation*}
$$

then

$$
\alpha_{n} \leq \frac{1}{2}\left(\beta_{n}+\sum_{2 \leq \ell \leq n-1} \ell \beta_{\ell} \beta_{n-\ell+1}\right)
$$

Since it is shown in [7] (p.22) that

$$
\sum_{2 \leq \ell \leq n-1} \frac{\ell(\ell-2)!(n-\ell-1)!}{(n-2)!} \leq \frac{5}{2} \text { for } n \geq 3
$$

we then obtain

$$
\alpha_{n} \leq \frac{1}{2}\left(1+\frac{5}{2 \Theta}\right) \beta_{n}
$$

Hence, it suffices to take

$$
\begin{equation*}
\Theta \geq \frac{5}{2} \tag{23}
\end{equation*}
$$

for having $\alpha_{n} \leq \beta_{n}$, which finally proves that

$$
\phi_{n} \leq \sqrt{m_{0}} K_{1}\left(\Theta K_{1}\right)^{n-2}(n!)^{\tau^{\prime}}(n-2)!, \quad n \geq 2
$$

provided that conditions $(20),(22),(23)$ on $\Theta$ and $K_{1}$ are satisfied. We can take for example

$$
\begin{equation*}
\Theta=8, K=8 K_{1}=\max \left\{\frac{9 \sqrt{m_{0}}}{\rho}, \frac{8 a c c_{01} m_{0}^{3 / 2}}{\rho^{2}}\right\} \tag{24}
\end{equation*}
$$

The first conclusion is that there exists $K>0$ depending only on $c, c_{01}, \rho, m_{0}, a$ such that

$$
\begin{equation*}
\phi_{n} \leq \sqrt{m_{0}} K^{n-1}(n!)^{1+\tau^{\prime}} \text { for } 1 \leq n \leq p \tag{25}
\end{equation*}
$$

Proof of Lemma 15. Estimate of $\sum_{n=1}^{p} \mathbf{\Phi}_{n}$.

First we have for $\left\|u_{0}\right\| \leq \delta$ and from Lemma 2.10 of $[7]$

$$
\begin{aligned}
\left\|\sum_{1 \leq n \leq p} \boldsymbol{\Phi}_{n}\left(u_{0}\right)\right\| & \leq \sum_{1 \leq n \leq p} \phi_{n} \delta^{n} \\
& \leq \sum_{1 \leq n \leq p} \frac{\sqrt{m_{0}}}{K}(\delta K)^{n}(n!)^{1+\tau^{\prime}} \\
& \leq \delta \sqrt{m_{0}} \sum_{1 \leq n \leq p}\left(\delta K p^{1+\tau^{\prime}}\right)^{n-1} .
\end{aligned}
$$

Let us choose $p$ such that

$$
\begin{equation*}
p=\left[\frac{1}{(2 \delta K)^{b}}\right], b=\frac{1}{1+\tau^{\prime}}, \tag{26}
\end{equation*}
$$

where [.] denotes the integer part of a number, then

$$
\left\|\sum_{1 \leq n \leq p} \boldsymbol{\Phi}_{n}\left(u_{0}\right)\right\| \leq 2 \delta \sqrt{m_{0}},
$$

and for $\delta<\rho /\left(2 \sqrt{m_{0}}\right)$ we have $\left\|u_{0}+\boldsymbol{\Phi}\left(u_{0}\right)\right\|<\rho$.
Proof of Lemma 16. Estimate of the remainder $\boldsymbol{\rho}\left(u_{0}\right)$.
We estimate each term $\mathfrak{R}_{k}$ separately.
Step 1. First we estimate $\mathfrak{R}_{1}\left(u_{0}\right)$. We have for every $\delta<\rho /\left(4 \sqrt{m_{0}}\right)$, and $p$ satisfying (26)

$$
\begin{align*}
\left\|\Re_{1}\left(u_{0}\right)\right\| \leq \sum_{q \geq p+1} c c_{01}\left(\frac{2 \delta \sqrt{m_{0}}}{\rho}\right)^{q} & \leq c c_{01}\left(\frac{2 \delta \sqrt{m_{0}}}{\rho}\right)^{p} \\
& \leq 2 c c_{01}\left(\frac{1}{2}\right)^{p+1}  \tag{27}\\
& \leq 2 c c_{01}\left(\frac{1}{2}\right)^{\frac{1}{(2 \delta K)^{b}}} \\
& \leq 2 c c_{01} e^{-\frac{\ln 2}{(2 \delta K)^{b}}} .
\end{align*}
$$

Step 2. For estimating $\mathfrak{R}_{2}\left(u_{0}\right)$ we have for $\delta<\rho /\left(4 \sqrt{m_{0}}\right)$

$$
\begin{aligned}
\left\|\mathfrak{R}_{2}\left(u_{0}\right)\right\| & \leq c c_{01} \sum_{2 \leq \ell \leq p, q \geq p+1} \ell \phi_{\ell} \sqrt{m_{0}} \delta^{\ell-1}\left(\frac{2 \delta \sqrt{m_{0}}}{\rho}\right)^{q} \\
& \leq c c_{01} \sum_{2 \leq \ell \leq p} \ell \phi_{\ell} \sqrt{m_{0}} \delta^{\ell-1}\left(\frac{2 \delta \sqrt{m_{0}}}{\rho}\right)^{p} .
\end{aligned}
$$

Now, for $p$ satisfying (26)

$$
\begin{aligned}
\sum_{2 \leq \ell \leq p} \ell \phi_{\ell} \delta^{\ell-1} & \leq \sum_{2 \leq n \leq p} \sqrt{m_{0}}(K \delta)^{n-1} n(n!)^{1+\tau^{\prime}} \\
& \leq \sum_{2 \leq n \leq p} \sqrt{m_{0}} p\left(K \delta p^{1+\tau^{\prime}}\right)^{n-1} \\
& \leq \sqrt{m_{0}} 2 K \delta p^{2+\tau^{\prime}} \leq \sqrt{m_{0}} p \leq \frac{\sqrt{m_{0}}}{(2 \delta K)^{b}}
\end{aligned}
$$

Hence, for $\delta<\delta_{1}=\min \left\{\rho /\left(4 \sqrt{m_{0}}\right), \frac{1}{2 K(2 e)^{1+\tau^{\prime}}}\right\}$ an using that that for $x \geq 2$, $\ln x \leq x \frac{\ln 2}{2}$, we get that

$$
\begin{align*}
\left\|\Re_{2}\left(u_{0}\right)\right\| & \leq \frac{c c_{01} m_{0}}{(2 \delta K)^{b}}\left(\frac{2 \delta \sqrt{m_{0}}}{\rho}\right)^{p} \\
& \leq \frac{2 c c_{01} m_{0}}{(2 \delta K)^{b}}\left(\frac{1}{2}\right)^{p+1} \\
& \leq \frac{2 c c_{01} m_{0}}{(2 K \delta)^{b}} e^{-\frac{\ln 2}{(2 \delta K)^{b}}} \\
& \leq 2 c c_{01} m_{0} e^{-\frac{\ln 2}{2(2 \delta K)^{b}}} \tag{28}
\end{align*}
$$

Step 3. We now estimate $\mathfrak{R}_{3}\left(u_{0}\right)$ :

$$
\left\|\Re_{3}\left(u_{0}\right)\right\| \leq c c_{01} \sum_{\substack{2 \leq q \leq p, 1 \leq k_{j} \leq p \\ p+1 \leq k_{1}+\cdots+k_{q}=n \leq q p}}\left(\frac{\sqrt{m_{0}}}{K \rho}\right)^{q}(\delta K)^{n}\left(k_{1}!\right)^{1+\tau^{\prime}} \ldots\left(k_{q}!\right)^{1+\tau^{\prime}}
$$

and from (24) we have $\frac{\sqrt{m_{0}}}{K \rho}=r \leq 1 / 9$ and from (26) we have $K \delta \leq \frac{1}{2 p^{1+\tau^{\prime}}}$. Hence,

$$
\begin{aligned}
\left\|\Re_{3}\left(u_{0}\right)\right\| & \leq c c_{01} \sum_{\substack{2 \leq q \leq p, 1 \leq k_{j} \leq p \\
p+1 \leq k_{1}+\cdots+k_{q}=n \leq q p}} r^{q} \frac{1}{2^{n}}\left(\frac{k_{1}!}{p^{k_{1}}}\right)^{1+\tau^{\prime}} \ldots\left(\frac{k_{q}!}{p^{k_{q}}}\right)^{1+\tau^{\prime}} \\
& \leq \frac{c c_{01}}{2^{p+1}} \sum_{2 \leq q \leq p} r^{q}\left(\sum_{1 \leq j \leq p}\left(\frac{j!}{p^{j}}\right)^{1+\tau^{\prime}}\right)^{q} .
\end{aligned}
$$

Moreover, we have

$$
\sum_{1 \leq j \leq p}\left(\frac{j!}{p^{j}}\right)^{1+\tau^{\prime}} \leq \frac{1}{p^{1+\tau^{\prime}}}+\sum_{2 \leq j \leq p} \frac{1}{p^{1+\tau^{\prime}}}=\frac{1}{p^{\tau^{\prime}}}
$$

hence, since $\frac{r}{p^{\tau^{\prime}}} \leq r \leq 1 / 9$

$$
\begin{align*}
\left\|\Re_{3}\left(u_{0}\right)\right\| & \leq \frac{c c_{01}}{2^{p+1}} \sum_{2 \leq q \leq p}\left(\frac{r}{p^{\tau^{\prime}}}\right)^{q} \\
& \leq \frac{c c_{01}}{72 \cdot 2^{p+1}} \leq \frac{c c_{01}}{72} e^{-\frac{\ln 2}{(2 \delta K)^{b}}} \tag{29}
\end{align*}
$$

Step 4. Finally, for the estimate of $\Re_{4}\left(u_{0}\right)$ we have by the same way

$$
\begin{align*}
& \left\|\Re_{4}\left(u_{0}\right)\right\| \leq c c_{01} \frac{m_{0} p^{1+\tau^{\prime}}}{2^{p+1}} \sum_{\substack{2 \leq \ell \leq p, \ell+k_{1}+\cdots+k_{q}=n+1 \leq p+2}} r^{q} \ell\left(\frac{\ell!}{p^{\ell}}\right)^{1+\tau^{\prime}}\left(\frac{k_{1}!}{p^{k_{1}}}\right)^{1+\tau^{\prime}} \ldots\left(\frac{k_{q}!}{p^{k_{q}}}\right)^{1+\tau^{\prime}} \\
& \leq c c_{01} \frac{m_{0} p^{1+\tau^{\prime}}}{2^{p+1}} \sum_{2 \leq q \leq p, 2 \leq \ell \leq p} r^{q} \ell\left(\frac{\ell!}{p^{\ell}}\right)^{1+\tau^{\prime}}\left(\sum_{1 \leq j \leq p}\left(\frac{j!}{p^{j}}\right)^{1+\tau^{\prime}}\right)^{q} \\
& \leq c c_{01} \frac{m_{0} p^{1+\tau^{\prime}}}{2^{p+1}} \sum_{2 \leq q \leq p}\left(\frac{r}{p^{\tau^{\prime}}}\right)^{q} \frac{p}{p^{\tau^{\prime}}} \\
& \leq \frac{c c_{01}}{72} \frac{m_{0} p^{2\left(1-\tau^{\prime}\right)}}{2^{p+1}} \\
& \leq \frac{c c_{01}}{72} m_{0} p^{2\left(1-\tau^{\prime}\right)} e^{-\frac{\ln 2}{(2 \delta K)^{b}}} \\
& \leq \frac{c c_{01} m_{0}}{72} e^{-\frac{\ln 2}{2(2 \delta K)^{b}}}, \tag{30}
\end{align*}
$$

provided that $\delta \leq \delta_{0}=\min \left\{\delta_{1}, \delta_{2}\right\}$ where $\delta_{2}$ is small enough, such that

$$
4\left(1-\tau^{\prime}\right) \ln \left(2 K \delta_{2}\right)^{-b} \leq\left(2 K \delta_{2}\right)^{-b} \ln 2
$$

this condition being empty for $\tau^{\prime} \geq 1$.
Now collecting estimates (27), (28), (29), (30), proves Theorem 1.

## 4 Analytic center manifolds up to Exponentially small terms

### 4.1 Finite dimensional case. Proof of theorem 5

This subsection is entirely devoted to the proof of Theorem 5 which ensures the existence of analytic center manifolds up to Exponentially small terms.

We notice that the diophantine condition (4) is automatically satisfied, since

$$
\left\langle\alpha, \lambda^{(0)}\right\rangle \in i \mathbb{R}
$$

and

$$
\left|\left\langle\alpha, \lambda^{(0)}\right\rangle-\lambda_{j}^{(1)}\right| \geq \gamma
$$

for all $\alpha \in \mathbb{N}^{m_{0}} \backslash\{0\}$ and all eigenvalues $\lambda_{j}^{(1)}$ of $\mathbf{L}_{1}$. Hence Theorem 1 applies directly, ensuring that there exists a polynomial $\boldsymbol{\Phi}: E_{0} \rightarrow E_{1}$ such that the change of variable in $E_{1}$

$$
u_{1}=v_{1}+\mathbf{\Phi}\left(u_{0}\right)
$$

transforms the system (1) into the following system in $E_{0} \times E_{1}$

$$
\begin{equation*}
\frac{d \widetilde{u}}{d t}=\mathbf{F}(\widetilde{u})+\widetilde{\boldsymbol{\rho}}(\widetilde{u}), \tag{31}
\end{equation*}
$$

where $\widetilde{u}=\left(u_{0}, v_{1}\right) \in E_{0} \times E_{1}$ and

$$
\mathbf{F}(\widetilde{u})=\binom{\mathbf{L}_{0} u_{0}+\mathbf{R}^{(0)}\left(u_{0}, v_{1}\right)}{\mathbf{L}_{1} v_{1}+\mathbf{R}^{(1)}\left(u_{0}, v_{1}\right)}, \widetilde{\boldsymbol{\rho}}(\widetilde{u})=\binom{0}{\boldsymbol{\rho}\left(u_{0}\right)}
$$

with

$$
\sup _{\left\|u_{0}\right\| \leq \delta}\left\|\boldsymbol{\rho}\left(u_{0}\right)\right\| \leq M e^{-\frac{w}{\delta}}
$$

For $\widetilde{\boldsymbol{\rho}} \equiv 0$, the truncated system

$$
\begin{equation*}
\frac{d \widetilde{u}}{d t}=\mathbf{F}(\widetilde{u}), \tag{32}
\end{equation*}
$$

admits the invariant manifold

$$
\widetilde{\mathcal{M}}_{0}^{\prime}=\left\{\widetilde{u}=\left(u_{0}, v_{1}\right) \in E_{0} \times E_{1} / v_{1}=0\right\}
$$

which appears to be an analytic center manifold (see [8], or [19] and references therein). In original coordinates this manifold reads

$$
\mathcal{M}_{0}^{\prime}=\left\{u \in \mathbb{R}^{m} / u=u_{0}+\boldsymbol{\Phi}\left(u_{0}\right)\right\}
$$

which is analytic since $\boldsymbol{\Phi}$ is polynomial.
Our aim is now to prove that for the full system (1), i.e. when $\widetilde{\boldsymbol{\rho}} \not \equiv 0$, this manifold is close to any center manifold up to an exponentially small term. For that purpose we see the full system in new coordinates (31) as a perturbation of the truncated system (32) by the exponentially small term $\widetilde{\boldsymbol{\rho}}(\widetilde{u})$.

We introduce three scalar parameters $(C, \varepsilon, \nu) \in[0,1]^{3}$ and consider the analytic vector field

$$
\begin{equation*}
\frac{d \widehat{v}}{d t}=\mathbf{V}(\widehat{v}, \varepsilon, C, \nu):=\mathbf{L} \widehat{v}+\frac{1}{\varepsilon} \mathbf{R}(\varepsilon \widehat{v})+\frac{\nu e^{\frac{C}{\varepsilon}}}{\varepsilon} \widetilde{\boldsymbol{\rho}}(\varepsilon \widehat{v}) . \tag{33}
\end{equation*}
$$

For $\nu=0$, (33) admits an analytic center manifold $\widetilde{\mathcal{M}_{0}^{\prime}}$ obtained from $\mathcal{M}_{0}^{\prime}$ by the scaling $\widetilde{u}=\varepsilon \widehat{v}$.

Since for every $(C, \varepsilon, \nu) \in[0,1]^{3}$ and every $\widehat{v} \in E_{0} \times E_{1}$ satisfying $\|\widehat{v}\| \leq$ $\delta_{0}=\frac{w}{2}$,

$$
\nu \frac{e^{\frac{C}{\varepsilon}}}{\varepsilon}\|\widetilde{\boldsymbol{\rho}}(\varepsilon \widehat{v})\| \leq \frac{\nu}{\varepsilon} e^{\frac{C-w / \delta_{0}}{\varepsilon}} \leq \frac{\nu}{\varepsilon} e^{-\frac{1}{\varepsilon}} \leq \mathfrak{m} \nu<\infty
$$

where $\mathfrak{m}:=\sup _{x \geq 0}\left(x e^{-x}\right)$, we know (see [19]) that there is a family of center manifolds $\mathcal{M}_{\varepsilon, \nu}^{x \geq 0}$ for $\|\widehat{v}\|+\varepsilon+|\nu|+C \leq r$ with $r<1$ which holds for

$$
0 \leq C \leq \frac{r}{2} \quad \text { and } \quad\|\widehat{v}\|+\varepsilon+|\nu| \leq \frac{r}{2} .
$$

Since we can choose $C_{*}<\min \left(\frac{r}{2}, \frac{4 w}{r}\right)$ such that for every $\nu \in[0, r / 4]$ and every $\widehat{v} \in E_{0} \times E_{1}$ and $\left.\left.\varepsilon \in\right] 0,1\right]$ satisfying $\|\widehat{v}\|+\varepsilon \leq \frac{r}{4}$,

$$
\frac{e^{\frac{C_{*}}{\varepsilon}}}{\varepsilon}\|\widetilde{\boldsymbol{\rho}}(\varepsilon \widehat{v})\| \leq \frac{1}{\varepsilon} e^{\frac{C_{*}-4 w / r}{\varepsilon}} \leq \frac{\mathfrak{m}}{4 w / r-C_{*}}<\infty
$$

the value $\nu=e^{-\frac{C_{*}}{\varepsilon}}$ is eligible for a center manifold which corresponds to the original system rescaled. The regularity results on center manifolds allow to claim that the graph satisfies

$$
u=u_{0}+\boldsymbol{\Phi}\left(u_{0}\right)+\boldsymbol{\Psi}\left(u_{0}\right),
$$

with constants $M$ and $C^{\prime}$ such that

$$
\left\|\boldsymbol{\Psi}\left(u_{0}\right)\right\| \leq M e^{-\frac{C^{\prime}}{\delta}}, \text { for }\left\|u_{0}\right\| \leq \delta
$$

Notice that we loose analyticity only in the term $\boldsymbol{\Psi}$ which is exponentially small.

### 4.2 Infinite dimensional case

The above result extends to the infinite dimensional case in the following way which needs an adapted assumption to replace Lemma 11. Indeed, still
in $\mathbb{R}^{m}$, and assuming that $\mathbf{L}_{0}$ is diagonal in $E_{0}$ where the norm is such that $e^{L_{0} t}$ is an isometry, we can solve the homological equation (14) in $\mathcal{H}_{n}$ in setting

$$
\mathbf{v}(t)=\Phi_{n}\left[\left(e^{\mathbf{L}_{0} t} u_{0}\right)^{(n)}\right] \in \mathcal{H}_{n}
$$

then

$$
\begin{equation*}
\frac{d \mathbf{v}(t)}{d t}=\mathbf{L}_{1} \mathbf{v}(t)+\mathbf{F}_{n}\left[\left(e^{\mathbf{L}_{0} t} u_{0}\right)^{(n)}\right] \tag{34}
\end{equation*}
$$

and it is easy to see that the unique solution which is allowed to possibly grow as $e^{\eta|t|}$ as $t \rightarrow \pm \infty$, with $\eta \in[0, \gamma[$, is given by
$\mathbf{v}(t)=\int_{-\infty}^{t} e^{\mathbf{L}_{1}^{-}(t-s)} \mathbf{P}_{-} \mathbf{F}_{n}\left[\left(e^{\mathbf{L}_{0} s} u_{0}\right)^{(n)}\right] d s-\int_{t}^{\infty} e^{\mathbf{L}_{1}^{+}(t-s)} \mathbf{P}_{+} \mathbf{F}_{n}\left[\left(e^{\mathbf{L}_{0} s} u_{0}\right)^{(n)}\right] d s$,
where the linear operators $\mathbf{P}_{ \pm}$are the projections commuting with $\mathbf{L}_{1}$, corresponding to the separation of its spectrum into eigenvalues with positive or negative real parts, and $\mathbf{L}_{1}^{ \pm}=\mathbf{P}_{ \pm} \mathbf{L}_{1}$. Moreover $\mathbf{v}(t)$ is smooth and bounded for $t \in \mathbb{R}$, and $t=0$ gives
$\Phi_{n}\left[\left(u_{0}\right)^{(n)}\right]=\int_{-\infty}^{0} e^{-\mathbf{L}_{1}^{-} s} \mathbf{P}_{-} \mathbf{F}_{n}\left[\left(e^{\mathbf{L}_{0} s} u_{0}\right)^{(n)}\right] d s-\int_{0}^{\infty} e^{-\mathbf{L}_{1}^{+} s} \mathbf{P}_{+} \mathbf{F}_{n}\left[\left(e^{\mathbf{L}_{0} s} u_{0}\right)^{(n)}\right] d s$,
and there is a constant $a$ depending only on the bounds of $e^{-\mathbf{L}_{1}^{-} s}$ for $s<0$ and of $e^{-\mathbf{L}_{1}^{+} s}$ for $s>0$ such that

$$
\phi_{n} \leq a\left|\mathbf{F}_{n}\right|_{2, n}
$$

Formula (35) which is valid in the finite dimensional space $E_{1}$ leads to a basic assumption for the center manifold theorem as formulated in [21], which is verified in many cases of physical interest (see examples in [21]).

Let $\mathcal{X}, \mathcal{Y}, \mathcal{Z}$ be (real or complex) Banach spaces such that

$$
\mathcal{Y} \hookrightarrow \mathcal{Z} \hookrightarrow \mathcal{X}
$$

with continuous embeddings. We consider a differential equation in $\mathcal{X}$, of the form

$$
\begin{equation*}
\frac{d u}{d t}=\mathbf{L} u+\mathbf{R}(u) \tag{37}
\end{equation*}
$$

in which we assume that the following holds.
$\mathrm{h}: 1$ Hypothesis 17 We assume that $\mathbf{L}$ and $\mathbf{R}$ in (37) have the following properties:
(a) $\mathbf{L} \in \mathcal{L}(\mathcal{Y}, \mathcal{X})$;
(b) There exists $\rho>0$ such that $\mathbf{R}: \mathcal{Y} \rightarrow \mathcal{Z}$ is analytic in the ball $\|u\|_{\mathcal{Y}} \leq$ $\rho$ and satisfies (2) and (3).

Besides the Hypothesis 17, we make two further assumptions on the linear operator $\mathbf{L}$, which are essential for the center manifold theorem.
$\mathrm{h}: 2$ Hypothesis 18 (Spectral decomposition) Consider the spectrum $\sigma$ of $\mathbf{L}$, and write

$$
\sigma=\sigma_{+} \cup \sigma_{0} \cup \sigma_{-}
$$

in which
$\sigma_{+}=\{\lambda \in \sigma ; \operatorname{Re} \lambda>0\}, \quad \sigma_{0}=\{\lambda \in \sigma ; \operatorname{Re} \lambda=0\}, \quad \sigma_{-}=\{\lambda \in \sigma ; \operatorname{Re} \lambda<0\}$.
We assume that
(a) there exists a positive constant $\gamma>0$ such that

$$
\inf _{\lambda \in \sigma_{+}}(\operatorname{Re} \lambda)>\gamma, \quad \sup _{\lambda \in \sigma_{-}}(\operatorname{Re} \lambda)<-\gamma ;
$$

(b) the set $\sigma_{0}$ consists of a finite number of eigenvalues with finite algebraic multiplicities and geometric multiplicity one.

This decomposition of the spectrum allows to define a projection $\mathbf{P}_{0} \in$ $\mathcal{L}\left(\mathcal{X}, E_{0}\right)$ on the finite-dimensional invariant "central" space $E_{0}$, which commutes with $\mathbf{L}$. The complementary projection $\mathbf{P}_{h}=\mathbb{I}-\mathbf{P}_{0}$ is also a projection commuting with $\mathbf{L}$, bounded in $\mathcal{X}_{h}=\mathbf{P}_{h} \mathcal{X}$ as well as in $\mathcal{Y}_{h}=\mathbf{P}_{h} \mathcal{Y}$ and $\mathcal{Z}_{h}=\mathbf{P}_{h} \mathcal{Z}$. The restriction of $\mathbf{L}$ to $\mathcal{Y}_{h}$ is denoted by $\mathbf{L}_{h}$.
$\mathrm{h}: 3$ Hypothesis 19 (Linear equation) For any $\eta \in[0, \gamma]$ and any

$$
f \in \mathcal{C}_{\eta}\left(\mathbb{R}, \mathcal{Z}_{h}\right)=\left\{v \in \mathcal{C}^{0}\left(\mathbb{R}, \mathcal{Z}_{h}\right) ;\left.\|v\|\right|_{\mathcal{C}_{\eta}}=\sup _{t \in \mathbb{R}}\left(e^{-\eta|t|}\|v(t)\| \mathcal{Z}_{h}\right)<\infty\right\},
$$

the linear problem

$$
\frac{d u_{h}}{d t}=\mathbf{L}_{h} u_{h}+f(t),
$$

has a unique solution $u_{h}=\mathbf{K}_{h} f \in \mathcal{C}_{\eta}\left(\mathbb{R}, \mathcal{Y}_{h}\right)$. Furthermore, the linear map $\mathbf{K}_{h}$ belongs to $\mathcal{L}\left(\mathcal{C}_{\eta}\left(\mathbb{R}, \mathcal{Z}_{h}\right), \mathcal{C}_{\eta}\left(\mathbb{R}, \mathcal{Y}_{h}\right)\right)$, and there exists a continuous map $C:[0, \gamma] \rightarrow \mathbb{R}$ such that

$$
\left\|\mathbf{K}_{h}\right\|_{\mathcal{L}\left(\mathcal{C}_{\eta}\left(\mathbb{R}, \mathcal{Z}_{h}\right), \mathcal{C}_{\eta}\left(\mathbb{R}, \mathcal{Y}_{h}\right)\right)} \leq C(\eta) .
$$

Then, we have the following theorem, which extends Theorem 5 to infinite dimensional cases:

## Theorem 20 (Center manifold analytic up to exp. small term)

Assume that the Hypotheses 17, 18, and 19 hold. Then for any $k>2$, there exists a polynomial $\mathbf{\Phi}: E_{0} \rightarrow E_{h}$ of degree $O(1 / \delta)$, with $\mathbf{\Phi}(0)=0$, $D \boldsymbol{\Phi}(0)=0$, a neighborhood $\mathcal{O}$ of 0 in $\mathcal{Y}$, and a map $\boldsymbol{\Psi} \in \mathcal{C}^{k}\left(E_{0}, \mathcal{Y}_{h}\right)$ which is $O\left(e^{-\frac{C}{\delta}}\right)$ for $\left\|u_{0}\right\|_{E_{0}} \leq \delta$ and a certain constant $C>0$, such that the manifold

$$
\begin{equation*}
\mathcal{M}_{0}=\left\{u_{0}+\boldsymbol{\Phi}\left(u_{0}\right)+\mathbf{\Psi}\left(u_{0}\right) ; u_{0} \in E_{0}\right\} \subset \mathcal{Y} \tag{38}
\end{equation*}
$$

has the following properties.
(a) $\mathcal{M}_{0}$ is locally invariant, i.e., if $u$ is a solution of (37) satisfying $u(0) \in$ $\mathcal{M}_{0} \cap \mathcal{O}$ and $u(t) \in \mathcal{O}$ for all $t \in[0, T]$, then $u(t) \in \mathcal{M}_{0}$ for all $t \in[0, T]$.
(b) $\mathcal{M}_{0}$ contains the set of bounded solutions of (37) staying in $\mathcal{O}$ for all $t \in \mathbb{R}$, i.e., if $u$ is a solution of (37) satisfying $u(t) \in \mathcal{O}$ for all $t \in \mathbb{R}$, then $u(0) \in \mathcal{M}_{0}$.

Proof. We use the result proved in [21], complemented by the proof of Theorem 5, for which we need to use Hypothesis 19 to solve the homological equation (14), as in (36), by

$$
\mathbf{\Phi}_{n}\left[u_{0}^{(n)}\right]=\left.\mathbf{K}_{h} \mathbf{F}_{n}\left[\left(e^{\mathbf{L}_{0} \cdot} u_{0}\right)^{(n)}\right]\right|_{t=0}
$$

and to obtain the basic estimate

$$
\phi_{n} \leq a\left|\mathbf{F}_{n}\right|_{2, n}
$$

## 5 Case of Elliptic vector fields

Consider now the system (1) in $\mathbb{R}^{n}$ when both spectra of $\mathbf{L}_{0}$ and $\mathbf{L}_{1}$ lie on the imaginary axis. This is the natural situation for nonlinear vibrating systems, typically with a large number of coupled nonlinear oscillators. This section is devoted to the proof of theorem 8 .

Theorem 1 applies and it results that the manifold $\mathcal{M}_{0}^{\prime}$ defined by

$$
u=u_{0}+\boldsymbol{\Phi}\left(u_{0}\right)
$$

which has the dimension of $E_{0}$ and is tangent to $E_{0}$ in 0 , is "nearly" invariant. More precisely, assume that the initial condition at $t=0$ is such that $\left.v_{1}\right|_{t=0}=0$, i.e. $\left.u\right|_{t=0} \in \mathcal{M}_{0}^{\prime}$. Then consider the second component of the vector field (6). If the remainder $\boldsymbol{\rho}\left(u_{0}\right)$ would be identically 0 , the manifold $\mathcal{M}_{0}^{\prime}$ would be an invariant manifold, since $v_{1}(t)=0$ would be the unique solution of the initial value problem. Now assume $v_{1}(0)=0$ and that $u_{0}(t)$ satisfies for $t \in[0, T]$

$$
\left\|u_{0}(t)\right\| \leq \delta .
$$

Then (6) and the estimate for $\mathbf{R}^{(1)}$ gives as soon as $\left\|v_{1}(t)\right\| \leq \delta$ for $t \in[0, T]$

$$
\left\|v_{1}(t)\right\| \leq c \delta \int_{0}^{t}\left\|e ^ { \mathbf { L } _ { 1 } ( t - s ) } \left|\left\|\mid v_{1}(s)\right\| d s+M t e^{-\frac{w}{\delta^{\hbar}}}\right.\right.
$$

For any $\xi>0$, there exists $C=\beta(\nu) \xi^{-(\nu-1)}$ where $\nu$ is the maximal index of eigenvalues of $\mathbf{L}_{1}$, such that for any $t \in \mathbb{R}$

$$
\left\|e^{\mathbf{L}_{1} t}\right\| \leq C e^{\xi|t|}
$$

then by Gronwall Lemma we get

$$
\left\|v_{1}(t)\right\| \leq M e^{-\frac{w}{\delta^{\delta}}}\left\{t+\frac{c C \delta}{(c C \delta+\xi)^{2}} e^{(c C \delta+\xi) t}\right\}
$$

and in choosing $\xi=(c \beta \delta)^{1 / \nu}$

$$
\begin{equation*}
\left\|v_{1}(t)\right\| \leq M e^{-\frac{w}{\delta^{b}}}\left\{t+\frac{1}{4(c \beta \delta)^{1 / \nu}} e^{2 t(c \beta \delta)^{1 / \nu}}\right\} \tag{39}
\end{equation*}
$$

which shows that $\left\|v_{1}(t)\right\|$ stays smaller than $M_{1} e^{-\frac{w}{2 \delta^{b}}}$ for $t=O\left(\delta^{-[b+1 / \nu]}\right)$. This means that the trajectory stays exponentially close to the manifold $\mathcal{M}_{0}^{\prime}$ for a very long time of order $O\left(\delta^{-[b+1 / \nu]}\right)$ and it achieves the proof of theorem 8 .

## A Norm of the inverse of the homological operator

Lemma 21 Let $\mathbf{L}$ be a linear operator in $\mathbb{C}^{m}$ and assume that the linear operator $\mathbf{L}$ is the direct sum of two linear operators $\mathbf{L}_{0}$ on $E_{0}$ (dim $\left.m_{0}\right)$, and $\mathbf{L}_{1}$ on $E_{1}\left(\operatorname{dim} m_{1}\right)$, such that $\mathbf{L}_{0}$ is diagonalizable with eigenvalues $\lambda_{1}^{(0)}, \cdots, \lambda_{m_{0}}^{(0)}$ and that there exist constants $0<\gamma \leq 1, \tau \geq 0$ such that $\Lambda_{\alpha, j}:=\left\langle\alpha, \lambda^{(0)}\right\rangle-\lambda_{j}^{(1)}$ satisfies

$$
\begin{equation*}
\left|\Lambda_{\alpha, j}\right| \geq \frac{\gamma}{|\alpha|^{\tau}} \tag{40}
\end{equation*}
$$

for any $\alpha \in \mathbb{N}^{m_{0}} \backslash\{0\}$, and any eigenvalue $\lambda_{j}^{(1)}$ of $\mathbf{L}_{1}$.
Let $\left(e_{k}\right)_{1 \leq k \leq m}$ be the canonical basis of $\mathbb{C}^{m}$. We assume that $\left(e_{k}\right)_{1 \leq k \leq m_{0}}$ is a basis of eigenvectors of $\boldsymbol{L}_{0}$ :

$$
\boldsymbol{L}_{0} e_{k}=\lambda_{j}^{(0)} e_{k} .
$$

Moreover we also assume that $f_{j}=e_{m_{0}+j}$ with $1 \leq j \leq m_{1}$, is a basis of generalized eigenvectors in which $\boldsymbol{L}_{1}$ is under Jordan complex normal form, i.e.

$$
\boldsymbol{L}_{1} f_{j}=\lambda_{j}^{(1)} f_{j}+\delta_{j-1} f_{j-1}
$$

where $\delta_{0}=0$ and where $\delta_{j}=0$ if $\lambda_{j}^{(1)} \neq \lambda_{j-1}^{(1)}$ and $\delta_{j}=0$ or 1 otherwise.
Let $\mathcal{H}$ be the set of all polynomials from $E_{0}$ to $E_{1}$ and let $\mathcal{H}_{n}$ be the subset of homogeneous polynomials of degree $n$. Finally let us denote by $\mathcal{A}: \mathcal{H} \rightarrow \mathcal{H}$ the homological operator defined by

$$
(\mathcal{A} \boldsymbol{\Phi})\left(u_{0}\right)=D \boldsymbol{\Phi}\left(u_{0}\right) \boldsymbol{L}_{0} u_{0}-\boldsymbol{L}_{1} \boldsymbol{\Phi}\left(u_{0}\right) .
$$

Then,
LemAcala (a) $\mathcal{A}$ maps $\mathcal{H}_{n}$ into $\mathcal{H}_{n}$ and the spectrum of its restriction to $\mathcal{H}_{n},\left.\mathcal{A}\right|_{\mathcal{H}_{n}}$, is given by

$$
\sigma\left(\left.\mathcal{A}\right|_{\mathcal{H}_{n}}\right):=\left\{\Lambda_{\alpha, j}=\left\langle\alpha, \lambda^{(0)}\right\rangle-\lambda_{j}^{(1)} / \alpha \in \mathbb{N}^{m_{0}},|\alpha|=n, 1 \leq j \leq m_{1}\right\} .
$$

LemAcalb (b) $\left.\mathcal{A}\right|_{\mathcal{H}_{n}}$ is invertible in the subspace $\mathcal{H}_{n}$ and

$$
\left.\left|\left||\mathcal{A}|_{\mathcal{H}_{n}}^{-1}\right|\right|\right|_{2}:=\left.\sup _{|\Phi|_{2}=1}|\mathcal{A}|_{\mathcal{H}_{n}}^{-1} \boldsymbol{\Phi}\right|_{2} \leq \nu \gamma^{-\nu} n^{\tau \nu} .
$$

Proof of (a). Let us denote by $\boldsymbol{P}_{\alpha, j}$ with $\alpha \in \mathbb{N}^{m_{0}},|\alpha|=n$ and $1 \leq j \leq m_{1}$ be the basis of $\mathcal{H}_{n}$ given by

$$
\boldsymbol{P}_{\alpha, j}\left(u_{0}\right)=\left(u_{0,1}\right)^{\alpha_{1}} \cdots\left(u_{0, m_{0}}\right)^{\alpha_{m_{0}}} f_{j}
$$

where $u_{0}=\sum_{k=1}^{m_{0}} u_{0, k} e_{k}$. Then we check that

$$
\begin{equation*}
\left.\mathcal{A}\right|_{\mathcal{H}_{n}} \boldsymbol{P}_{\alpha, j}=\Lambda_{\alpha, j} \boldsymbol{P}_{\alpha, j}-\delta_{j-1} \boldsymbol{P}_{\alpha, j-1} \tag{41}
\end{equation*}
$$

Let us order this basis by lexicographical order, i.e. $\boldsymbol{P}_{\alpha, j}<\boldsymbol{P}_{\beta, \ell}$ if the first non zero integer $\beta_{1}-\alpha_{1}, \cdots, \beta_{m_{0}}-\alpha_{m_{0}}, \ell-j$ is positive. Within this
order, the matrix $M_{\left.\mathcal{A}\right|_{\mathcal{H}_{n}}}$ of $\left.\mathcal{A}\right|_{\mathcal{H}_{n}}$ in the basis $\boldsymbol{P}_{\alpha, j}$ is upper triangular. More precisely, it is the direct sum of $m_{1} \times m_{1}$ matrices $M_{\alpha}$,

$$
M_{\left.\mathcal{A}\right|_{\mathcal{H}_{n}}} \bigoplus_{\alpha \in \mathbb{N}^{m_{0},|\alpha|=m_{0}}} M_{\alpha}, \quad M_{\alpha}=\left(\begin{array}{cccc}
\Lambda_{\alpha, 1} & \delta_{1} & 0 & 0  \tag{42}\\
0 & \ddots & \ddots & 0 \\
0 & 0 & \ddots & \delta_{m_{1}} \\
0 & 0 & 0 & \Lambda_{\alpha, m_{1}}
\end{array}\right)
$$

Hence the spectrum of $\left.\mathcal{A}\right|_{\mathcal{H}_{n}}$ is given by

$$
\sigma\left(\left.\mathcal{A}\right|_{\mathcal{H}_{n}}\right):=\left\{\Lambda_{\alpha, j}=\left\langle\alpha, \lambda^{(0)}\right\rangle-\lambda_{j}^{(1)} / \alpha \in \mathbb{N}^{m_{0}},|\alpha|=n, 1 \leq j \leq m_{1}\right\}
$$

Proof of (b). Since by hypothesis, for every $\alpha \in \mathbb{N}^{m_{0}}$ and every $1 \leq j \leq m_{1}$, $\left|\Lambda_{\alpha, j}\right| \geq \frac{\gamma}{\left.\alpha\right|^{\tau}}>0,\left.\mathcal{A}\right|_{\mathcal{H}_{n}}$ is invertible and (42) ensures that

$$
M_{\left.\mathcal{A}\right|_{\mathcal{H}} ^{-1}}^{-1}=M_{\left.\mathcal{A}\right|_{\mathcal{H}}}^{-1}=\bigoplus_{\alpha \in \mathbb{N}^{m_{0}},|\alpha|=m_{0}} M_{\alpha}^{-1}
$$

Moreover $M_{\alpha}$ is block diagonal

$$
M_{\alpha}=\bigoplus_{r=1}^{q} B_{j_{r}, p_{r}} \quad \text { with } \quad B_{j, p}=\left(\begin{array}{cccc}
\Lambda_{\alpha, j} & 1 & 0 & 0 \\
0 & \ddots & \ddots & 0 \\
0 & 0 & \ddots & 1 \\
0 & 0 & 0 & \Lambda_{\alpha, j+p}
\end{array}\right)
$$

where $1 \leq j_{r} \leq m_{1}$ and $0 \leq p_{r} \leq \nu$ where $\nu$ is the maximal index of the eigenvalues of $\boldsymbol{L}_{1}$. For a polynomial $\boldsymbol{\Phi} \in \mathcal{H}_{n}$, we can write

$$
\mathbf{\Phi}=\sum_{|\alpha|=n} \sum_{j=1}^{m_{1}} \Phi_{\alpha, j} \boldsymbol{P}_{\alpha, j}=\sum_{|\alpha|=n} \sum_{r=1}^{q} \sum_{j=j_{r}}^{j_{r}+p_{r}} \Phi_{\alpha, j} \boldsymbol{P}_{\alpha, j}
$$

Then for $\boldsymbol{\Psi} \in \mathcal{H}_{n}, \mathcal{A} \boldsymbol{\Phi}=\boldsymbol{\Psi}$ if and only if, for every $\alpha \in \mathbb{N}^{m_{0}}$ with $|\alpha|=n$ and every $1 \leq r \leq q$

$$
\left\{\begin{array}{cl}
\Phi_{\alpha, j_{r}+p_{r}} & =\Lambda_{\alpha, j_{r}+p_{r}}^{-1} \Psi_{\alpha, j_{r}+p_{r}} \\
\Phi_{\alpha, j_{r}+p_{r}-1} & =\Lambda_{\alpha, j_{r}+p_{r}-1}^{-1} \Psi_{\alpha, j_{r}+p_{r}-1}-\left(\Lambda_{\alpha, j_{r}+p_{r}} \Lambda_{\alpha, j_{r}+p_{r}-1}\right)^{-1} \Psi_{\alpha, j_{r}+p_{r}} \\
\vdots & \\
\Phi_{\alpha, j_{r}}= & \Lambda_{\alpha, j_{r}}^{-1} \Psi_{\alpha, j_{r}+p_{r}-1}-\left(\Lambda_{\alpha, j_{r}} \Lambda_{\alpha, j_{r}+1}\right)^{-1} \Psi_{\alpha, j_{r}+1} \\
& +\cdots+(-1)^{p_{r}-1}\left(\Lambda_{\alpha, j_{r}} \cdots \Lambda_{\alpha, j_{r}+p_{r}}\right)^{-1} \Psi_{\alpha, j_{r}+p_{r}}
\end{array}\right.
$$

Then observe that for every $\alpha \in \mathbb{N}^{m_{0}}$ with $|\alpha|=n$,

$$
\max _{\substack{1 \leq r \leq q \\ \leq j \leq \ell \leq j_{r}+p_{r}}}\left(\Lambda_{\alpha, j} \cdots \Lambda_{\alpha, \ell}\right)^{-1} \leq \gamma^{-\nu} n^{\nu \tau}
$$

Thus, for every $\alpha \in \mathbb{N}^{m_{0}}$ with $|\alpha|=n$, every $1 \leq r \leq q$ and every $j_{r} \leq j \leq$ $j_{r}+p_{r}$,

$$
\left|\Phi_{\alpha, j}\right| \leq \gamma^{-\nu} n^{\tau \nu}\left(\left|\Psi_{j_{r}}\right|+\cdots+\left|\Psi_{j_{r}+p_{r}}\right|\right.
$$

Hence, since $\left\langle P_{\alpha, j}, P_{\beta, \ell}\right\rangle_{\mathcal{H}}=0$ for $(j, \alpha) \neq(\ell, \beta)$ and since $\left|P_{j, \alpha}\right|_{2}=\left|P_{\ell, \alpha}\right|_{2}=$ $\alpha$ !, we have

$$
\begin{aligned}
|\boldsymbol{\Phi}|_{2}^{2} & =\left|\mathcal{A}^{-1} \boldsymbol{\Psi}\right|_{2}^{2} \\
& =\sum_{|\alpha|=n} \sum_{r=1}^{q} \sum_{j=j_{r}}^{j_{r}+p_{r}}\left|\Phi_{\alpha, j}\right|^{2}\left|\boldsymbol{P}_{\alpha, j}\right|_{2}^{2} \\
& \leq\left(\gamma^{-\nu} n^{\nu \tau}\right)^{2} \sum_{|\alpha|=n} \sum_{r=1}^{q} \sum_{j=j_{r}}^{j_{r}+p_{r}}\left(\sum_{\ell=j_{r}}^{j_{r}+p_{r}}\left|\Psi_{\alpha, \ell}\right|\right)^{2}\left|\boldsymbol{P}_{\alpha, j}\right|_{2}^{2} \\
& \leq \nu\left(\gamma^{-\nu} n^{\nu \tau}\right)^{2} \sum_{|\alpha|=n} \sum_{r=1}^{q} \sum_{j=j_{r}}^{j_{r}+p_{r}} \sum_{\ell=j_{r}}^{j_{r}+p_{r}}\left|\Psi_{\alpha, \ell}\right|^{2}\left|\boldsymbol{P}_{\alpha, j}\right|_{2}^{2} \\
& \leq \nu^{2}\left(\gamma^{-\nu} n^{\nu \tau}\right)^{2} \sum_{|\alpha|=n} \sum_{r=1}^{q} \sum_{\ell=j_{r}}^{j_{r}+p_{r}}\left|\Psi_{\alpha, \ell}\right|^{2}\left|\boldsymbol{P}_{\alpha, \ell}\right|_{2}^{2} \\
& =\left(\nu \gamma^{-\nu} n^{\nu \tau}\right)^{2}|\boldsymbol{\Psi}|_{2}^{2}
\end{aligned}
$$

Hence, $\left|\mathcal{A}^{-1} \boldsymbol{\Psi}\right|_{2} \leq \nu \gamma^{-\nu} n^{\nu \tau}|\Psi|_{2}$.

## References

Car [1] J. Carr. Applications of center manifold theory. Applied Mathematical Sciences, 35. Springer-Verlag, New York-Berlin, 1981.
[2] Delshams A. \& Gutierrez P. Effective stability and KAM theory. J. Differential Equations 128, no.2, 415-490 (1996).

Fasso [3] Fassò F. Lie series method for vector fields and Hamiltonian perturbation theory. J. Appl. Math. Phys. 41 843-864 (1990).
[4] Giorgilli A. \& Posilicano A. Estimates for normal forms of differential equation near an equilibrium point. $Z A M P$, vol. 39, 713-732 (1988).

Gior-Mu [5] A.Giorgilli, D.Muraro. Exponentially stable manifolds in the neighborhood of elliptic equilibria. Boll. Unione Mat. Ital. Sez. B9 (2006), 1-20.

Gro-Schn [6] M.Groves, G.Schneider. Modulating pulse solutions to quadratic quasilinear wave equations over exponentially long length scales. Com. Math. Phys. 278 (2008), 567-625.

Io-Lomb [7] G.Iooss, E.Lombardi. Polynomial normal forms with exponentially small remainder for analytic vector fields. J. Diff. Equ. 212 (2005) 1-61.

Kelley [8] A.Kelley. The stable, center-stable, center, center-unstable, unstable manifolds. J.Differential equations 3 (1967), 546-570.

## Jezequel

[9] L. Jézéquel and C.H. Lamarque, Analysis of non-linear dynamical systems by the normal form theory, Journal of Sound and Vibration 149 (1991) (3), pp. 429459.

Lom [10] E.Lombardi. Oscillatory integrals and phenomena beyond all algebraic orders, with applications to homoclinic orbits in reversible systems. Lecture Notes in Maths vol. 1741. Springer 2000.

NekI [11] Nekhoroshev, N.N. An exponential estimate of the time of stability of nearly integrable Hamiltonian systems, I. Usp. Mat. Nauk 32 5-66 (1997); Russ. Math. Surv. 32, 1-65 (1977).

NekII [12] Nekhoroshev, N.N. An exponential estimate of the time of stability of nearly integrable Hamiltonian systems, II. Tr. Semin. Petrovsk. 5., 5-50 (1979); In:Oleineik, O.A.(ed.) Topics in Modern Mathematics, Petrovskii Semin., no. 5. New York: Consultant Bureau (1985).

Mars-McCr [13] J.Marsden, M.McCracken. The Hopf bifurcation and its applications. Appl. Math. Sci. 19, Springer, New York, 1976.

Poschel [14] Pöschel J. Nekoroshev estimates for quasi convex hamiltonian systems. Math. Z. 213, 187-216, (1993) .

Saw [15] S.W. Shaw and C. Pierre, Non-linear normal modes and invariant manifolds, J. Sound Vib. 150 (1) (1991), pp. 170173.

Sij [16] J.Sijbrand. Studies in nonlinear stability and bifurcation theory. PhD Thesis, University of Utrecht, 1981.

Touz [17] C.Touzé, O.Thomas, A.Chaigne. Hardening/softening behavior in nonlinear oscillations of structural systems using non-linear normal modes. J. Sound and vibration 273 (2004) 77-101.

Touz2 [18] C. Touzé, M. Amabili : Non-linear normal modes for damped geometrically non-linear systems: application to reduced-order modeling of harmonically forced structures, Journal of Sound and Vibration, vol. 298, No. 4-5, pp. 958-981, 2006.

V [19] A. Vanderbauwhede. Centre manifolds, normal forms and elementary bifurcations. Dynamics reported, Vol. 2, 89-169, Dynam. Report. Ser. Dynam.Systems Appl., 2, Wiley, Chichester, 1989.

VStr [20] S.Van Strien. Center manifolds are not $\mathcal{C}^{\infty}$. Math. Z. 166 (1979), 143145.

V-I [21] A. Vanderbauwhede, G. Iooss. Center manifold theory in infinite dimensions. Dynamics reported: expositions in dynamical systems, 125163, Dynam. Report. Expositions Dynam. Systems (N.S.), 1, Springer, Berlin, 1992.

