# Renormalisation of flows on the multidimensional torus close to a $K T$ frequency vector 

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#### Abstract

We use a renormalisation operator $\mathcal{R}$ acting on a space of vector fields on $\mathbb{T}^{d}, d \geq 2$, to prove the existence of a submanifold of vector fields equivalent to constant. The result comes from the existence of a fixed point $\omega$ of $\mathcal{R}$ which is hyperbolic. This is done for a certain class $K T_{d}$ of frequency vectors $\omega \in \mathbb{R}^{d}$, called of Koch type. The transformation $\mathcal{R}$ is constructed using a time rescaling, a linear change of basis plus a periodic non-linear map isotopic to the identity, which we derive by a "homotopy trick".


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## 1 Introduction

In this paper we consider a renormalisation transformation $\mathcal{R}$ that acts on a Banach space of analytic vector fields $X$ on $\mathbb{T}^{d}=\mathbb{R}^{d} / \mathbb{Z}^{d}, d \geq 2$. Its domain is an open ball around a non-zero constant vector field $\omega \in \mathbb{R}^{d}$ (that generates a linear flow). The operator $\mathcal{R}$ is a $C^{1}$ infinite-dimensional dynamical system with a hyperbolic fixed point at $\omega$. We show that the stable and unstable manifolds, given by the local dynamics of $\mathcal{R}$ around $\omega$, yield different equivalence classes of vector fields. The elements of the codimension- $(d-1)$ stable manifold are equivalent to $\omega$, as $\mathcal{R}$ is designed to be an equivalence between flows. In other words, the renormalisation asymptotically contracts locally any $X$ into a ( $d-1$ )-parameter family of constant vector fields that includes $\omega$. The map $\mathcal{R}$ basically consists of the composition of a time rescaling with a diffeomorphism of $\mathbb{T}^{d}$ isotopic to a linear automorphism in $G L(d, \mathbb{Z})$.

This idea of renormalising vector fields is due to MacKay [8], and we follow an approach inspired by the work of Koch [5] on renormalisation for $d$-degrees of freedom analytic Hamiltonian systems. The latter applies to the problem of stability of invariant tori associated with a frequency vector $\omega$.

The linear transformation is the main feature of the operator. It is a change of the basis of $\mathbb{T}^{d}$ lifted to $\mathbb{R}^{d}$, made to enlarge the region around the orbits of the linear flow. By iterating this, we look closer at chosen regions, but periodicity and the whole torus are kept at each stage.

We also need to rescale the time as orbits take longer to cross the "new" torus. The topological characteristics of the trajectories are not affected by such transformation.

The non-linear part of $\mathcal{R}$, isotopic to the identity, is obtained by a homotopy method to reduce the perturbation. We use a flow of coordinate changes (diffeomorphisms) to find one which fully eliminates some of the Fourier terms of $X$. Those are chosen to be the easiest to do so, since we keep the "resonant" terms associated with "small denominators" (as it appears in the usual KAM theory). We are not concerned with trying to cancel all the perturbation terms because the linear transformation shifts resonant to "non-resonant" terms. Eventually all perturbation terms are eliminated while iterating $\mathcal{R}$. The separation between resonant and non-resonant terms appears in [5].

If a vector field is a fixed point of the above procedure, it means that its orbits exhibit self-similarity between the different scales. This is trivially deduced for the fixed point $\omega$.

The existence of the linear part, based on an arithmetical condition [5], depends on certain conditions imposed on $\omega$. In particular, for $d=2$, that corresponds to the set of vectors with a "quadratic irrational" slope. Specifically for the two-dimensional case, it is defined in [7] a family of renormali-
sation iterative schemes allowing a full Lebesgue measure set of diophantine vectors. The methods involved therein are not easily generalisable to higher dimensions, it is required a suitable choice of a multidimensional continued fraction expansion algorithm.

This paper is organised as follows. In Section 2 we recall the basic ingredients of equivalence of flows, and in Section 3 we introduce the renormalisation idea through an example. We rigorously construct the operator $\mathcal{R}$ in Sections $4,5,6$ and 7 , with proofs of the statements also given in Section 9. Finally, in Section 8 we present the main result using the spectral properties of the derivative at the fixed point, with further discussion in Section 10.

## 2 Equivalence of flows

Consider a continuous vector field $X$ on $\mathbb{T}^{d}$ that generates the flow $\phi_{t}, \dot{\theta}=$ $X(\theta), \theta \in \mathbb{T}^{d}$, with lift $\Phi_{t}$ to the universal cover $\mathbb{R}^{d}$. Choosing a norm $\|\cdot\|$ in $\mathbb{R}^{d}$, we define $w_{X}\left(\theta_{0}\right)=\lim _{t \rightarrow \infty} \Phi_{t}\left(\theta_{0}\right) /\left\|\Phi_{t}\left(\theta_{0}\right)\right\|$ to be the winding ratio of $X$ for the orbit of $\theta_{0} \in \mathbb{T}^{d}$, if the limit exists and $\lim _{t \rightarrow \infty}\left\|\Phi_{t}\left(\theta_{0}\right)\right\|=\infty$. Otherwise, if $\Phi_{t}\left(\theta_{0}\right)$ is bounded we put $w_{X}\left(\theta_{0}\right)=0$. Also, if the limit does not exist or if $\left\|\Phi_{t}\left(\theta_{0}\right)\right\|$ is unbounded but does not tend to infinity, the winding ratio is not defined.

We say that two flows $\phi_{t}, \psi_{t}: \mathbb{T}^{d} \rightarrow \mathbb{T}^{d}$ are $C^{r}$-equivalent if there is a $C^{r}$ diffeomorphism $h: \mathbb{T}^{d} \rightarrow \mathbb{T}^{d}$ taking orbits of $\phi_{t}$ onto those of $\psi_{t}$, preserving orientation. We are allowing to have $h\left(\phi_{t}(\theta)\right)=\psi_{\tau(h(\theta), t)}(h(\theta))$, where $\tau(\theta, \cdot)$ is a homeomorphism of $\mathbb{R}$ for any $\theta \in \mathbb{T}^{d}$. This is the same to say that two vector fields $X, Y$ on $\mathbb{T}^{d}$ are equivalent if $\tau^{\prime} X \circ h=D h Y$ where $h$ is a flow equivalence as above and $\tau^{\prime}$ the time derivative of $\tau$. We emphasise that this relaxation of the usual requirement that $t$ be preserved ( $h \circ \phi_{t}=$ $\left.\psi_{t} \circ h\right)$ provides more satisfactory equivalence classes for flows, since we are in fact mainly interested in qualitative (topological) properties of the flow. Every $C^{r}$-equivalence is isotopic to a map with a lift to $\mathbb{R}^{d}$ in the group of the linear automorphisms of the lattice $\mathbb{Z}^{d}$ with determinant $\pm 1, G L(d, \mathbb{Z})$. The isotopy is given by periodic homeomorphisms on $\mathbb{T}^{d}$. These kind of coordinate changes preserve the $d$-torus structure and volume.

The set of winding ratios of a flow on $\mathbb{T}^{d}$ generated by $X$ is $w_{X}=$ $\left\{w_{X}(\theta): \theta \in \mathbb{T}^{d}\right\}$, which is called the winding set. Any automorphism $T: \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}$ induces the map $\hat{T}: \mathbb{S}^{d-1} \rightarrow \mathbb{S}^{d-1}$ given by $x \mapsto T x /\|T x\|$, where $\mathbb{S}^{d-1}=\left\{x \in \mathbb{R}^{d}:\|x\|=1\right\}$. Let $T$ be in $G L(d, \mathbb{Z})$. Then the winding set is invariant up to the action of $\hat{T}$, with a $C^{r}$-equivalence $h$ isotopic to a map with lift $T$. That is, if $X^{\prime}=(D h)^{-1} X \circ h$, then $w_{X^{\prime}}\left(\theta^{\prime}\right)=\hat{T} w_{X}(\theta)$, where $\theta^{\prime}=h^{-1}(\theta)$. In particular, the winding set is preserved by transformations isotopic to the identity.


Figure 1: Action of $T$ and $U$ on an orbit of the lifted flow of $X$.

## 3 Motivating example for $d=2$

We start by motivating the renormalisation procedure with a simple example. Consider the case $d=2$ and the linear flow described by the differential equation $\dot{\theta}=\omega$, with $\omega \in \mathbb{R}^{2}$ and $\theta \in \mathbb{T}^{2}$. Given an initial condition $\theta_{0} \in \mathbb{T}^{2}$ for $t=0$, the solution of the flow is $\phi_{t}\left(\theta_{0}\right)=\theta_{0}+\omega t \bmod 1, t>0$. The motion is simply a rotation with frequency vector $\omega$ and winding ratio $w_{\omega}=\omega /\|\omega\|$. If, for all non-zero integer vectors $k, k \cdot \omega:=k_{1} \omega_{1}+k_{2} \omega_{2} \neq 0$ (i.e. the slope $\omega_{2} / \omega_{1}$ is irrational, assuming $\omega_{1}$ does not vanish), then all the orbits are dense in the torus (the flow is minimal). Otherwise, they are closed curves (periodic orbits).

Considering a perturbation of a constant "irrational" vector field $\omega$ lifted to $\mathbb{R}^{2}$, we want to determine under which conditions there is still equivalence to $\omega$. Let $X$ be a vector field close to $\omega$ arising from a time-independent perturbation. We choose e.g. $\omega=(1, \gamma)$ where $\gamma=\frac{1+\sqrt{5}}{2}$ is the golden ratio. The main idea is to perform a change of basis, from the canonical base to $\{(0,1),(1,1)\}$, enlarging the region of $\mathbb{R}^{2}$ around the orbits of $\omega$ (see Figure 1). This is achieved by the linear transformation: $\theta^{\prime}=T^{-1} \theta$, where $T=\left[\begin{array}{ll}0 & 1 \\ 1 & 1\end{array}\right] \in G L(2, \mathbb{Z})$. The eigenvalues and eigenvectors of $T$ are given by: $T \omega=\gamma \omega$ and $T \Omega=-\frac{1}{\gamma} \Omega$, where $\Omega \perp \omega$. We also rescale the time, $t^{\prime}=\frac{1}{\gamma} t$, because the orbits take longer to cross the new torus. This transformation does not affect the unperturbed vector field $\omega$, as it is given by $\gamma D T^{-1} \circ T\left(\theta^{\prime}\right) \omega \circ T\left(\theta^{\prime}\right)=\gamma T^{-1} \omega=\omega$.

We should now consider a coordinate change $h\left(\theta^{\prime}\right)=U_{X} \circ T\left(\theta^{\prime}\right)$ with $U$ satisfying $U_{\omega}=\mathrm{Id}$. The fundamental requirement is that it has to cancel the growth of the wiggles in the orbits of $X$, which are enlarged by the use of $T$ (see orbit in Figure 1). The idea is to perform a non-linear, close to
the identity periodic coordinate transformation $U_{X}$, in order to remove as many perturbation terms as possible.

There are some terms which are more relevant - "resonant". To understand what they are and how they appear, consider the vector field $X$ in the form of its Fourier decomposition:

$$
\dot{\theta}=\omega+\varepsilon \sum_{k \in \mathbb{Z}^{2}} f_{k} e^{2 \pi i k \cdot \theta}
$$

with $f_{k} \in \mathbb{C}^{2}$ and $\varepsilon>0$ "small". We can approximately solve the above system using the unperturbed solution plus an order $\varepsilon$ term. So, on the universal cover,

$$
\begin{equation*}
\theta(t)=\theta_{0}+\omega t+\varepsilon f_{0} t+\varepsilon \sum_{k \in \mathbb{Z}^{2} \backslash\{0\}} \frac{f_{k}}{2 \pi i k \cdot \omega} e^{2 \pi i k \cdot\left(\theta_{0}+\omega t\right)}+\mathcal{O}\left(\varepsilon^{2}\right) \tag{1}
\end{equation*}
$$

for $t>0$. The resonant terms are the ones whose index $k$ is almost perpendicular to $\omega$.

We need to add some conditions to $\omega$ in order to avoid having small denominators in the solution of the flow. This is done in the usual KAMtype proofs by imposing a diophantine condition on $\omega$, i.e. the denominators $|k \cdot \omega|$ in (1) admit a lower bound. We will show in the following that it is enough to find $U$ such that it eliminates only "far from resonance terms". This is so because of the extra linear change of coordinates $T$ described before, that is responsible for a "shift" of resonant into non-resonant modes. In some cases, by iterating this process the orbits can be straightened.

We define $U$ for any non-zero vector $\omega$. Naturally, a restriction on $\omega$ will appear again in the procedure with $T$ (see also [7]), since this type of results fails for some vectors with irrational slope, in particular for some non-diophantine vectors [1].

## 4 Space of Analytic Vector Fields

The following is valid for any $d$-dimensional torus $\mathbb{T}^{d}, d \geq 2$, for a class of frequencies $\omega$ to be defined. The vector fields considered are inside a ball around $\omega$ in some adequate space, and can be regarded as maps of $\mathbb{R}^{d}$ by lifting their domains. We make use of the analyticity to extend to the complex domain, so we deal with complex analytic vector fields. We construct the renormalisation operator $\mathcal{R}$ and we look at the spectral properties of its derivative at the fixed point $\omega$, to relate to the local dynamics.

The use of analytic function spaces as the domain of the renormalisation functional operator $\mathcal{R}$ is justified by the usefulness of $\mathcal{R}$ being $C^{1}$. Otherwise the "picture" of $\mathcal{R}$ being a dynamical system with stable and unstable manifolds of the fixed points would vanish. The problem lies in the compositions that appear in $\mathcal{R}$, thus reducing the degree of differentiability of
the image of the renormalisation operator. Considering $C^{\infty}$ functions would also be a possibility, but complicate the technical parts of the method.

Let $r>0$ and the domain lifted to a complex neighbourhood of $\mathbb{R}^{d}$ :

$$
\mathcal{D}(r)=\left\{\theta \in \mathbb{C}^{d}:\|\operatorname{Im} \theta\|<\frac{r}{2 \pi}\right\}
$$

where $\|\cdot\|$ is the $\ell_{1}$-norm on $\mathbb{C}^{d}$. That is, $\|z\|=\sum_{i=1}^{d}\left|z_{i}\right|$, with $|\cdot|$ the usual norm on $\mathbb{C}$. We will also be using the inner product: $z \cdot z^{\prime}=\sum_{i=1}^{d} z_{i} \bar{z}_{i}^{\prime}$.

An analytic function $f: \mathcal{D}(r) \rightarrow \mathbb{C}^{d}, 2 \pi$-periodic in each variable $\theta_{i}$, is represented in Fourier series as

$$
f(\theta)=\sum_{k \in \mathbb{Z}^{d}} f_{k} e^{2 \pi i k \cdot \theta}
$$

with coefficients $f_{k} \in \mathbb{C}^{d}$. The Banach spaces $\left(\mathcal{A}_{d}(r),\|\cdot\|_{r}\right)$ and $\left(\mathcal{A}_{d}^{\prime}(r),\|\cdot\|_{r}^{\prime}\right)$ are the subsets of the set of these functions such that the respective norms

$$
\|f\|_{r}=\sum_{k \in \mathbb{Z}^{d}}\left\|f_{k}\right\| e^{r\|k\|} \quad \text { and } \quad\|f\|_{r}^{\prime}=\sum_{k \in \mathbb{Z}^{d}}(1+2 \pi\|k\|)\left\|f_{k}\right\| e^{r\|k\|}
$$

are finite.
Let $\rho>0$ be a fixed value in the following. Consider the vector fields of the form $X(\cdot)=\omega+f(\cdot)$, where $\omega \in \mathbb{R}^{d}$ and $f \in \mathcal{A}_{d}^{\prime}(\rho)$. Assume that $X$ has no equilibrium points, $X(\theta) \neq 0$ for any $\theta \in \mathbb{T}^{d}$. A condition like $\|f\|_{\rho}^{\prime}<\|\omega\|$ is enough to assure that there are no equilibria and, for any $\theta \in \mathbb{T}^{d}$, there is no $\theta^{\prime} \in \mathbb{T}^{d}$ such that $X\left(\theta^{\prime}\right) /\left\|X\left(\theta^{\prime}\right)\right\|=-X(\theta) /\|X(\theta)\|$. So, we could not have two symmetric winding ratios for such a vector field, since it would have implied having at least two symmetric normalised tangent vectors. The set of such vector fields on $\mathbb{T}^{2}$ corresponds to the class of the Poincaré flows (fixed-point-free with only one possible winding ratio, see e.g. [2]). We study these systems in more detail in [7].

## 5 Elimination of the Far from Resonance Terms

We will be dealing with the vector fields considered above decomposed in Fourier series: $X(\theta)=\sum_{k \in \mathbb{Z}^{d}} X_{k} e^{2 \pi i k \cdot \theta}$. As in [5], it is important to distinguish different classes of terms according to their respective indices.

Definition 5.1 For a fixed value of $\sigma>0$ and $\psi \in \mathbb{C}^{d}$, we define the far from resonance terms with respect to $\psi$ to be the ones whose indices are in

$$
I_{\sigma}^{-}(\psi)=\left\{k \in \mathbb{Z}^{d}:|\psi \cdot k|>\sigma\|k\|\right\} .
$$

Similarly, the resonant terms are $I_{\sigma}^{+}(\psi)=\mathbb{Z}^{d} \backslash I_{\sigma}^{-}(\psi)$. It is also useful to define the projections $\mathbb{I}_{\sigma}^{+}(\psi)$ and $\mathbb{I}_{\sigma}^{-}(\psi)$ for any vector field $X$ by

$$
\left[\mathbb{I}_{\sigma}^{ \pm}(\psi)\right] X(\theta)=\sum_{k \in I_{\sigma}^{ \pm}(\psi)} X_{k} e^{2 \pi i k \cdot \theta}
$$

and $\mathbb{I}=\mathbb{I}_{\sigma}^{+}(\psi)+\mathbb{I}_{\sigma}^{-}(\psi)$ is the identity operator.
The existence of a nonlinear change of coordinates $U$, close to the identity, that eliminates the far from resonance terms $I_{\sigma}^{-}(\omega)$ of a vector field $X$ in a neighbourhood of a non-zero constant vector field, is given by

Theorem 5.2 Let $0<\rho^{\prime}<\rho, \omega \in \mathbb{C}^{d} \backslash\{0\}$ and $0<\sigma<\|\omega\|$. If $X$ is in the open ball $\hat{B} \subset \mathcal{A}_{d}^{\prime}(\rho)$ centred at $\omega$ with radius

$$
\hat{\varepsilon}=\frac{\sqrt{6}-2}{12} \sigma \min \left\{\frac{\rho-\rho^{\prime}}{4 \pi}, \frac{3-\sqrt{6}}{6} \frac{\sigma}{\|\omega\|}\right\}
$$

there is an invertible coordinate change $U: \mathcal{D}\left(\rho^{\prime}\right) \rightarrow \mathcal{D}(\rho)$ close to the identity, satisfying

$$
\left[\mathbb{I}_{\sigma}^{-}(\omega)\right](D U)^{-1} X \circ U=0, \quad \text { and } \quad U=\operatorname{Id} \text { if }\left[\mathbb{I}_{\sigma}^{-}(\omega)\right] X=0
$$

The $\operatorname{map} \mathcal{U}: \hat{B} \rightarrow\left[\mathbb{I}_{\sigma}^{+}(\omega)\right] \mathcal{A}_{d}\left(\rho^{\prime}\right)$ given by $X \mapsto(D U)^{-1} X \circ U$ is analytic, and the derivative at every constant vector field $\omega+\psi$ in $\hat{B}$ is equal to $\mathbb{I}_{\sigma}^{+}(\omega)$. Moreover,

$$
\|\mathcal{U}(X)-\omega\|_{\rho^{\prime}} \leq 2\left(1+\max \left\{\frac{2}{3}(3-\sqrt{6}), 6(\sqrt{6}+2) \frac{\|\omega\|}{\sigma}\right\}\right)\|X-\omega\|_{\rho}^{\prime}
$$

In Section 9 we include a proof of the above theorem using the "homotopy method". It is worth noting here that since we are requiring to eliminate only the far from resonance terms outside $I_{\sigma}^{+}(\omega)$, we avoid the problem in the proof of dealing with small denominators (see Section 9.3).

## 6 Change of Basis

All the coordinate changes of $\mathbb{T}^{d}$ are isotopic to a linear transformation in $G L(d, \mathbb{Z})$ which preserves the structure and volume of the $d$-torus. That can be seen as a map on the frequency vectors space, acting on a vector field by shifting terms in the space of indices $\mathbb{Z}^{d}$ and changing coefficients.

Definition 6.1 For a fixed non-zero vector $\omega \in \mathbb{R}^{d}$ with rationally independent components, i.e. $k \cdot \omega \neq 0$ for any non-zero integer vector $k$, assume that $\omega=\left(1, \omega_{2}, \ldots, \omega_{d}\right)$ where its components are real algebraic numbers (roots
of non-zero polynomials over $\mathbb{Q}$ ). The vectors $c \omega, c \in \mathbb{R} \backslash\{0\}$, are said to be of Koch type, i.e. to belong to the set $K T_{d}$ if the algebraic extension of $\mathbb{Q}$ by the numbers $\omega_{2}, \ldots, \omega_{d}$ (the smallest field containing $\omega_{2}, \ldots, \omega_{d}$ and $\mathbb{Q}$ ) is of degree $d$.

An important property, the existence of a linear change of basis $T$ corresponding to the above class of vectors, follows from

Lemma 6.2 (Koch [5] - Lemma 4.1) A vector $\omega$ is in $K T_{d}$ if and only if there is an integral $d \times d$ matrix $T$, such that $0<\left|\lambda_{d}\right| \leq \cdots \leq\left|\lambda_{2}\right|<$ $1<\left|\lambda_{1}\right|$ and $T \omega=\lambda_{1} \omega$, where $\lambda_{1}, \ldots, \lambda_{d}$ are simple eigenvalues of $T$. Furthermore, some of the (integral) matrices satisfying these properties are in $G L(d, \mathbb{Z})$.

Remark 6.3 For $d=2$, the numbers $\omega_{2}$ that produce a degree 2 algebraic extension of $\mathbb{Q}$ are the quadratic irrationals (roots of non-zero degree 2 polynomials over $\mathbb{Q}$ ). These are characterised by an eventually periodic continued fraction expansion [6]. For $d=3$, the base for the algebraic extension of $\mathbb{Q}$ by $\alpha$ and $\beta$ is $\left\{1, \alpha, \alpha^{2}\right\}$, where $\alpha$ has to be a cubic irrational (root of a non-zero degree 3 polynomial over $\mathbb{Q}$ ). Then $\beta=c_{1}+c_{2} \alpha+c_{3} \alpha^{2}$, $c_{i} \in \mathbb{Q}$.

In what follows, a vector $\omega \in K T_{d}$ and a corresponding matrix $T \in$ $G L(d, \mathbb{Z})$ are chosen according to Lemma 6.2. Also, $\bar{\omega} \in \mathbb{R}^{d}$ is chosen to be the (unstable) $\lambda_{1}$-eigenvector of $T^{*}$ such that $\bar{\omega} \cdot \omega=1$. Notice that $K T_{d}$ is a subset of the diophantine vectors $D C(\beta)$ with $\beta=-1-\ln \left|\lambda_{1}\right| / \ln \left|\lambda_{2}\right|([5]-$ Corollary 4.2), i.e. there is a constant $C>0$ that satisfies $|\omega \cdot k|>C\|k\|^{-\beta-1}$ for every integer vector $k \neq 0$. In addition, $\prod_{j=1}^{d}\left|\lambda_{j}\right|=1$, where $\lambda_{j}$ are the eigenvalues of $T$, with $\lambda_{1} \in \mathbb{R}$ since it is the only one outside the unit circle. Also, $\bar{\omega} \cdot \omega^{(j)}=0$, with $\omega^{(j)}$ being the eigenvector of $T$ corresponding to $\lambda_{j}, 2 \leq j \leq d$, because $\bar{\omega} \cdot \omega^{(j)}=T^{*} \bar{\omega} \cdot T^{-1} \omega^{(j)}=\left(\lambda_{1} / \bar{\lambda}_{j}\right) \bar{\omega} \cdot \omega^{(j)}$. We denote the transpose matrix of $T$ by $T^{*}$ and will always assume that $\max \left\{\left\|T^{*} y\right\|:\|y\|=1, \omega \cdot y=0\right\}<\kappa$ for some $\kappa<1$.

From now on, we simplify the notation writing $I_{\sigma}^{ \pm}$and $\mathbb{I}_{\sigma}^{ \pm}$instead of the sets $I_{\sigma}^{ \pm}(\omega)$ and the projections $\mathbb{I}_{\sigma}^{ \pm}(\omega)$, respectively.

The next proposition determines that the change of basis corresponding to the matrix $T$ above, acting on the space of "resonant vector fields", is analyticity improving.

Proposition 6.4 Let $0<\kappa<1$ and $0<\rho^{\prime}<\rho$. If $\kappa \rho<\rho^{\prime}$, then, for some $0<\sigma<\frac{1}{2}\|\bar{\omega}\|^{-1}$, any $X \in \mathbb{I}_{\sigma}^{+} \mathcal{A}_{d}\left(\rho^{\prime}\right)$ has an analytic extension on $T \mathcal{D}(\rho)$. The linear map $\mathcal{T}(X)=T^{-1} X \circ T$ from $\mathbb{I}_{\sigma}^{+} \mathcal{A}_{d}\left(\rho^{\prime}\right)$ to $\mathcal{A}_{d}^{\prime}(\rho)$ is compact.

The proof is given in Section 9.4. The values of $\sigma$ have to be chosen such that

$$
I_{\sigma}^{+} \subset I^{\kappa}:=\left\{x \in \mathbb{R}^{d}:\left\|T^{*} x\right\| \leq \kappa\|x\|\right\}
$$

for a given $\kappa<1$. The proof of this proposition only requires that the resonant terms are inside the set $I^{\kappa}$. Thus, we could have different definitions for $I_{.}^{+}$. The different choices are restricted by the part of the proof of Theorem 5.2 in Section 9.3, that requires the denominator terms $|\omega \cdot k|$ to have at least a $\|k\|$-linear lower bound on $I^{-}$.

It is also important to remark that one assumes that the orthogonal hypersurface in $\mathbb{R}^{d}$ with respect to $\omega$ is contained in $I^{\kappa}$. This condition is indeed verified by choosing appropriate norms on $\mathbb{R}^{d}$ or by considering a sufficiently large power of $T$ instead. The indices of resonant terms are therefore shifted towards $I_{\sigma}^{-}$, eventually becoming far from resonance terms.

## 7 Renormalisation Operator

We are now in condition to construct a renormalisation operator $\mathcal{R}$ based on the coordinate transformations introduced before.

Fix $\omega \in K T_{d}, T$ and $\lambda_{j}, 1 \leq j \leq d$, as in Lemma 6.2. We denote by $\mathbb{E}(X)$ the average on $\mathbb{T}^{d}$ of a vector field $X$, i.e.

$$
\mathbb{E}(X)=\int_{\mathbb{T}^{d}} X(\theta) d \theta
$$

where $d \theta$ is the normalised Lebesgue measure on $\mathbb{T}^{d}$.
Definition 7.1 Consider the maps $\mathcal{U}$ and $\mathcal{T}$ given by Theorem 5.2 and Proposition 6.4, respectively. The renormalisation operator $\mathcal{R}$ is defined to be

$$
\mathcal{R}(X)=\frac{1}{\bar{\omega} \cdot \mathbb{E}(\tilde{X})} \tilde{X}, \quad \tilde{X}=\mathcal{T} \circ \mathcal{U}(X)
$$

for vector fields $X$.
The rescaling of time is chosen to preserve the term $\omega$, because without it the iterates of $\omega,(\mathcal{T} \circ \mathcal{U})^{n}(\omega)=\lambda_{1}^{-n} \omega, n \geq 0$, would go to zero. We could fix that by simply considering the product by $\lambda_{1}$ inside the calculation of $\tilde{X}$, but then all the elements of the one-parameter family $X_{\mu}=(1+\mu) \omega$, $\mu \in \mathbb{C}$, would be fixed points of the renormalisation. This "neutral" direction is contracted by the operator $\mathcal{R}$ given in the definition, allowing $\omega$ to be an isolated fixed point.

This choice of rescaling has the disadvantage of reducing the domain of the operator (see Proposition 7.2 below), when comparing with the size of the ball $\hat{B}$ obtained in Theorem 5.2. In any case, this can be solved by using the rescaling suggested above, thus working with a one-parameter family of fixed points and a larger domain.

Note also that, by Proposition $6.4, \sigma<\frac{1}{2}\|\omega\|^{-1}\|\omega\|_{2}^{2} \leq \frac{1}{2}\|\omega\|$, where $\|\cdot\|_{2}$ represents the Euclidean norm. This means that the radius of an open


Figure 2: Steps of the renormalisation operator $\mathcal{R}$ for $d=2$, on the lattice $\mathbb{Z}^{2}$, the indices of the Fourier coefficients (notice that each coefficient is a vector in $\left.\mathbb{C}^{d}\right)$.
ball given by Theorem 5.2 can be taken to be $\mathcal{C} \sigma^{2}\|\omega\|^{-1}$, assuming at least $\mathcal{C}<4$. Hence, the domain of $\mathcal{R}$ only contains equilibria-free vector fields.

The change of basis given by $T$ turns some of the resonant terms into $I_{\sigma}^{-}$where they will be completely eliminated by the operator $\mathcal{U}$ of the next iteration of $\mathcal{R}$. The use of such $\mathcal{U}$ is not essential, it would be enough to "sufficiently" reduce the resonant terms in a way that analyticity would not be lost. As the elimination is complete, we get a faster convergence to the fixed point and a much simpler and clearer analysis starting from the fact that the derivative of $\mathcal{U}$ is straightforward. If we would only apply $\mathcal{U}$, all the neutral directions that are eliminated by the use of $\mathcal{T}$ would persist (no elimination of the resonant terms), and the analyticity domain would not be fully recovered. The right combination of the two coordinate changes is responsible for the usefulness of $\mathcal{R}$ (see Figure 2 for the $d=2$ case).

Proposition 7.2 Let $\rho>0$. The renormalisation operator $\mathcal{R}$ is a welldefined analytic map from an open neighbourhood $B \subset \mathcal{A}_{d}^{\prime}(\rho)$ of $\omega$, to $\mathcal{A}_{d}^{\prime}(\rho)$.
Proof: Consider the complex-valued continuous functional $F(\tilde{X})=\bar{\omega}$. $\mathbb{E}(\tilde{X})$, and restrict its domain to a neighbourhood $\tilde{B} \subset \mathcal{A}_{d}^{\prime}(\rho)$ of $\lambda_{1}^{-1} \omega$ such that $F(\tilde{B}) \subset\left\{z \in \mathbb{C}:\left|\lambda_{1} z-1\right|<1 / 2\right\}$, i.e. $F$ is bounded away from zero in $\tilde{B}$.

The domain $B$ of $\mathcal{R}$ is a subset of $\hat{B}$ (the domain of $\mathcal{U}$ ) such that $\mathcal{T} \circ$ $\mathcal{U}(B) \subset \tilde{B}$. That is satisfied for a small enough choice of the radius of $B \subset \hat{B}$, since $\mathcal{U}(X)=\omega+\mathbb{I}_{\sigma}^{+} f+\mathcal{O}\left(\|f\|_{\rho}^{\prime 2}\right)$ and $\left|\lambda_{1} F(\mathcal{T} \circ \mathcal{U}(X))-1\right|=$ $\left|\lambda_{1}\left(\bar{\omega} \cdot T^{-1} \mathbb{E} f\right)+\mathcal{O}\left(\|f\|_{\rho}^{2}\right)\right|$, with $X=\omega+f$. It is then sufficient to have $\left|\bar{\omega} \cdot \mathbb{E} f+\mathcal{O}\left(\|f\|_{\rho}^{\prime 2}\right)\right|<\frac{1}{2}$. Notice that $\bar{\omega} \cdot T^{-1} \mathbb{E} f=T^{-1 *} \bar{\omega} \cdot \mathbb{E} f=\lambda_{1}^{-1}(\bar{\omega} \cdot \mathbb{E} f)$. Hence, as $F$ is bounded and analytic in $\tilde{B}$, the above, Theorem 5.2 and Proposition 6.4 prove the claim.

## 8 Hyperbolicity of the Fixed Point $\omega$

The transformation $\mathcal{R}$ was constructed in a way to hold an isolated fixed point at $\omega$. The local behaviour of $\mathcal{R}$ around it is given by the derivative, with which we hope to characterise the vector fields close enough to $\omega$, realising the existence or not of an analytic equivalence.

We can rewrite $\mathcal{R}$ by the expression $\mathcal{R}=\mathcal{F} \circ \mathcal{T} \circ \mathcal{U}$, where the timerescaling step is given by $\mathcal{F}(X)=[\bar{\omega} \cdot \mathbb{E}(X)]^{-1} X$ and $D \mathcal{F}\left(\lambda_{1}^{-1} \omega\right) f=\lambda_{1}^{-1}[f-$ $\bar{\omega} \cdot \mathbb{E}(f)] \omega$. The derivative of $\mathcal{R}$ at $\omega$ is then the linear map

$$
D \mathcal{R}(\omega) f=\mathcal{L} f-[\bar{\omega} \cdot \mathbb{E}(\mathcal{L} f)] \omega, \quad f \in \mathcal{A}_{d}^{\prime}(\rho), \quad \text { with } \quad \mathcal{L} f=\lambda_{1} \mathcal{T} \circ \mathbb{I}_{\sigma}^{+} f
$$

This operator is compact by Proposition 6.4 and using the fact that $\mathbb{I}_{\sigma}^{+}$is bounded. The eigenvalues of $\mathcal{L}$ are zero and those of $T^{-1}$ multiplied by $\lambda_{1}$. The vector $\omega$ gives the subspace corresponding to the eigenvalue 1. Since the renormalisation should not depend on the size of $\omega$, this direction is fully eliminated in $D \mathcal{R}(\omega)$ by the rescaling of time: $D \mathcal{R}(\omega) \omega=\omega-(\bar{\omega}$. $\omega) \omega=0$. There are $(d-1)$ unstable directions of $D \mathcal{R}(\omega)$ given by the other eigenvectors of $T^{-1}$, with eigenvalues $\lambda_{1} \lambda_{j}^{-1}, 2 \leq j \leq d$, of modulus greater than one (see Section 10).

The remaining eigenvalue of $D \mathcal{R}(\omega)$ is zero, meaning that the elimination of only the far from resonant terms in conjuction with the change of basis are sufficient to eliminate all the non-constant Fourier terms of a vector field close to $\omega$. Therefore, there is a codimension- $(d-1)$ manifold inside a neighbourhood of $\omega$ in $\mathcal{A}_{d}^{\prime}(\rho)$ being mapped by $\mathcal{R}$ into itself, and we have proved

Theorem 8.1 If $\omega \in K T_{d}$, then the constant vector field $\omega$ is a hyperbolic fixed point of $\mathcal{R}$ with a local codimension- $(d-1)$ stable manifold $\mathcal{W}^{s}(\omega)$ and a local $(d-1)$-dimensional unstable manifold $\mathcal{W}^{u}(\omega)$.

We have determined an analytic equivalence between the elements of $\mathcal{W}^{s}(\omega)$ and in particular with $\omega$, that is given by the operator $\lim _{n \rightarrow+\infty} \mathcal{R}^{n}$. The local stable manifold $\mathcal{W}^{s}(\omega)$ is the set of all $X$ in some neighbourhood $B$ with winding set $w_{X}=\omega /\|\omega\|$ (see Section 10 ). On the other hand, the local unstable manifold $\mathcal{W}^{u}(\omega)$ is the affine space of the vector fields in the form $X=\omega+v \in B$, where $v \in \operatorname{span}\left\{\omega^{(2)}, \ldots, \omega^{(d)}\right\}$, as this spectral subspace is $\mathcal{R}$-invariant. A schematic representation of the renormalisation scheme described is in Figure 3.

## 9 Proofs of Theorem 5.2 and Proposition 6.4

### 9.1 Preliminaries

We include here some technical details useful for the proofs.


Figure 3: Schematic renormalisation picture given by the action of $\mathcal{R}$ in a neighbourhood of $\omega$. Note that $\operatorname{dim} \mathcal{W}^{u}(\omega)=d-1$ whilst $\mathcal{W}^{s}(\omega)$ is an infinite dimensional manifold (of codimension $d-1$ ).

Firstly, note that by writing $f=\left(f^{1}, \ldots, f^{d}\right)$, one can rewrite the formulas of the norms on the spaces $\mathcal{A}_{d}(r)$ and $\mathcal{A}_{d}^{\prime}(r)$. In fact,

$$
\|f\|_{r}=\sum_{i=1}^{d}\left\|f^{i}\right\|_{r}, \quad f^{i} \in \mathcal{A}_{1}(r) \quad \text { and } \quad\|f\|_{r}^{\prime}=\sum_{i=1}^{d}\left\|f^{i}\right\|_{r}^{\prime}, \quad f^{i} \in \mathcal{A}_{1}^{\prime}(r),
$$

where

$$
\left\|f^{i}\right\|_{r}^{\prime}=\left\|f^{i}\right\|_{r}+\sum_{j=1}^{d}\left\|\partial_{j} f^{i}\right\|_{r}, \quad i=1, \ldots, d .
$$

One has $\|f\|_{r} \leq\|f\|_{r}^{\prime}, \mathcal{A}_{d}^{\prime}(r) \subset \mathcal{A}_{d}(r)$ and $\|f\|_{r^{\prime}} \leq\|f\|_{r}$, for $0<r^{\prime}<r$.
Define the inclusion maps $\mathcal{I}: \mathcal{A}_{d}(r) \rightarrow \mathcal{A}_{d}\left(r^{\prime}\right)$ and $\mathcal{I}^{\prime}: \mathcal{A}_{d}^{\prime}(r) \rightarrow \mathcal{A}_{d}^{\prime}\left(r^{\prime}\right)$, $0<r^{\prime}<r$, by the restriction of the domain $\mathcal{D}(r)$ to $\mathcal{D}\left(r^{\prime}\right)$, i.e. $\left.f \mapsto f\right|_{\mathcal{D}\left(r^{\prime}\right)}$. It is easy to check that $\mathcal{I}\left(\mathcal{A}_{d}(r)\right) \subset \mathcal{A}_{d}\left(r^{\prime}\right)$ and $\mathcal{I}^{\prime}\left(\mathcal{A}_{d}^{\prime}(r)\right) \subset \mathcal{A}_{d}^{\prime}\left(r^{\prime}\right)$. The inclusion maps $\mathcal{I}$ and $\mathcal{I}^{\prime}$ are compact linear operators.

Consider a bounded linear map $A: \mathcal{A}_{d}(r) \rightarrow \mathcal{A}_{d}(r)$ with $(A f)(\theta)=$ $A(\theta) f(\theta)$, where $A(\theta)$ is a linear map on $\mathbb{C}^{d}$. Its operator norm satisfies $\|A\| \leq \sum_{i, j=1}^{d}\left\|a_{i j}\right\|_{r}$, where $a_{i j} \in \mathcal{A}_{1}(r)$ such that, for each $\theta \in \mathcal{D}(r)$, $a_{i j}(\theta)$ are the coefficients of the matrix which represents $A(\theta)$ in the canonical basis. In particular, if $A$ is the derivative of a function $f \in \mathcal{A}_{d}(r)$, $(D f g)(\theta)=D f(\theta) g(\theta), g \in \mathcal{A}_{d}(r)$ and $\theta \in \mathcal{D}(r)$, then $\|D f\| \leq\|f\|_{r}^{\prime}$.

If $\alpha \in \mathcal{A}_{1}(r)$ and $f, g \in \mathcal{A}_{d}(r)$, then $\alpha f \in \mathcal{A}_{d}(r), f \cdot g \in \mathcal{A}_{1}(r),\|\alpha f\|_{r} \leq$ $\|\alpha\|_{r}\|f\|_{r}$ and $\|f \cdot g\|_{r} \leq\|f\|_{r}\|g\|_{r}$. Similarly, if $\alpha \in \mathcal{A}_{1}^{\prime}(r)$ and $f, g \in \mathcal{A}_{d}^{\prime}(r)$, then $\alpha f \in \mathcal{A}_{d}^{\prime}(r), f \cdot g \in \mathcal{A}_{1}^{\prime}(r),\|\alpha f\|_{r}^{\prime} \leq\|\alpha\|_{r}^{\prime}\|f\|_{r}^{\prime}$ and $\|f \cdot g\|_{r}^{\prime} \leq\|f\|_{r}^{\prime}\|g\|_{r}^{\prime}$. In this way, the spaces $\mathcal{A}_{d}(r)$ and $\mathcal{A}_{d}^{\prime}(r)$ are Banach algebras.

Lemma 9.1 Let $0<r^{\prime}<r$ and $f \in \mathcal{A}_{d}^{\prime}(r)$. If $u \in \mathcal{A}_{d}\left(r^{\prime}\right)$ and $\|u\|_{r^{\prime}}<$ $\left(r-r^{\prime}\right) / 4 \pi$, then

1. $\|f(\operatorname{Id}+u)\|_{r^{\prime}} \leq\|f\|_{\left(r+r^{\prime}\right) / 2}$,
2. $\|D f(\operatorname{Id}+u)\| \leq\|f\|_{\left(r+r^{\prime}\right) / 2}^{\prime}$,
3. $\|f(\operatorname{Id}+u)-f\|_{r^{\prime}} \leq\|f\|_{\left(r+r^{\prime}\right) / 2}^{\prime}\|u\|_{r^{\prime}}$,
4. $\|D f(\operatorname{Id}+u)-D f\| \leq \frac{4 \pi}{r-r^{\prime}}\|f\|_{r}^{\prime}\|u\|_{r^{\prime}}$.

Proof: Knowing that $\left\|e^{2 \pi i k \cdot u}\right\|_{r^{\prime}} \leq e^{2 \pi\|k\|\|u\|_{r^{\prime}}}$ one gets 1 and 2. The mean value theorem gives 3 . To prove 4 apply again the mean value theorem now to $\left\|\partial_{j} f^{i}(\operatorname{Id}+u)-\partial_{j} f^{i}\right\|_{r^{\prime}}$ and also

$$
\begin{aligned}
\left\|\partial_{j} f^{i}\right\|_{\left(r+r^{\prime}\right) / 2}^{\prime} & =\sum_{k \in \mathbb{Z}^{d}} 2 \pi\left|k_{j}\right|(1+\|k\|)\left|f_{k}^{i}\right| e^{r\|k\|} e^{-\left(r-r^{\prime}\right)\|k\| / 2} \\
& \leq \frac{4 \pi}{r-r^{\prime}} \sum_{k \in \mathbb{Z}^{d}}(1+\|k\|)\left|f_{k}^{i}\right| e^{r\|k\|}=\frac{4 \pi}{r-r^{\prime}}\left\|f^{i}\right\|_{r}^{\prime}
\end{aligned}
$$

where we have used the inequality: $\sup _{t \geq 0} t e^{-\beta t} \leq 1 / \beta$ for $\beta>0$.

### 9.2 Homotopy Method

The sometimes called "homotopy trick" has been used in different problems, such as by Moser to show that all smooth volume forms on a compact orientable manifold are equivalent up to a diffeomorphism ([4] - Section 5.1e). Other examples of its application are proofs of the Darboux theorem or the Poincaré lemma ([4] - Sections 5.5.9 and A.3.11, respectively), and Roussarie's proof of Morse's lemma. This procedure is also related to the "deformation method" used in different setups in KAM theory. It consists of a flow of symplectomorphisms that reduce in each iterate the size of the perturbation of an integrable Hamiltonian [9, 3].

In the following fix $\sigma>0$, thus dropping the index $\sigma$ for the sets $I^{ \pm}$and corresponding projections $\mathbb{I}^{ \pm}$. Choose also $\rho$ and $\rho^{\prime}$ such that $0<\rho^{\prime}<\rho$. Let $0<\hat{\delta}<\frac{1}{2}$ and $\hat{\varepsilon}$ be a positive constant which will be determined along the proof and contains the restrictions on the size of the perturbation of $\omega$ depending on $\sigma$ and the norm of $\omega$, as will be seen in Section 9.3.

For vector fields in the form $X=\omega+f$, consider the open neighbourhood $\mathcal{E}$ of the term $f$ :

$$
\mathcal{E}=\left\{f \in \mathcal{A}_{d}^{\prime}(\rho):\|f\|_{\rho}^{\prime}<\hat{\varepsilon}\right\} .
$$

The coordinate transformation $U$ is written as $U=\mathrm{Id}+u$, with Id as the identity transformation, and $u$ in the open ball $\mathcal{B}$ of radius $\hat{\delta}$ in $\mathbb{I}^{-} \mathcal{A}_{d}^{\prime}\left(\rho^{\prime}\right)$, i.e.

$$
\mathcal{B}=\left\{u \in \mathbb{I}^{-} \mathcal{A}_{d}^{\prime}\left(\rho^{\prime}\right):\|u\|_{\rho^{\prime}}^{\prime}<\hat{\delta}\right\} .
$$

For a given function $f \in \mathcal{E}$, consider the operator $F: \mathcal{B} \rightarrow \mathbb{I}^{-} \mathcal{A}_{d}\left(\rho^{\prime}\right)$,

$$
F(u)=\mathbb{I}^{-}(I+D u)^{-1}[\omega+f(\mathrm{Id}+u)],
$$

where $I$ is the identity operator. This is simply the transformed vector field $\mathbb{I}^{-}(D U)^{-1} X \circ U$, with $X=\omega+f$.

Lemma 9.2 The derivative of $F$ at $u \in \mathcal{B}$ is the linear map from $\mathbb{I}^{-} \mathcal{A}_{d}^{\prime}\left(\rho^{\prime}\right)$ to $\mathbb{I}^{-} \mathcal{A}_{d}\left(\rho^{\prime}\right)$,

$$
\begin{aligned}
D F(u) h= & \mathbb{I}^{-}(I+D u)^{-1}[D f \circ(\operatorname{Id}+u) h \\
& \left.-D h(I+D u)^{-1}(\omega+f \circ(\operatorname{Id}+u))\right],
\end{aligned}
$$

with $h$ chosen such that $u+h \in \mathcal{B}$.
Proof: We need to compute the linear term on $h$ of $F(u+h)-F(u)$. As $u \in \mathcal{B}$, we have that the bounded linear operator $D u$, from $\mathcal{A}_{d}\left(\rho^{\prime}\right)$ into itself, satisfies $\|D u\| \leq\|u\|_{\rho^{\prime}}^{\prime}<1$. Hence, $(I+D u)^{-1}=\sum_{n \geq 0}(-D u)^{n}$. Using the following formulas:

$$
\begin{aligned}
& (D u+D h)^{n}=D u^{n}+\sum_{j=0}^{n-1} D u^{j} D h D u^{n-1-j}+\mathcal{O}\left(\|h\|_{\rho^{\prime}}^{\prime}\right), \quad n \geq 1 \\
& \sum_{n=1}^{+\infty} \sum_{j=0}^{n-1}(-1)^{n} D u^{j} D h D u^{n-1-j}=-(I+D u)^{-1} D h(I+D u)^{-1}
\end{aligned}
$$

and the Taylor expansion of $f$ around $\theta+u(\theta), \theta \in \mathcal{D}(\rho)$, that gives

$$
f \circ(\operatorname{Id}+u+h)=f \circ(\operatorname{Id}+u)+D f \circ(\operatorname{Id}+u) h+\mathcal{O}\left(\|h\|_{\rho^{\prime}}^{\prime}{ }^{2}\right),
$$

we get

$$
\begin{aligned}
F(u+h)-F(u)= & \mathbb{I}^{-} \sum_{n=0}^{+\infty}(-1)^{n}\left[(D u+D h)^{n}(\omega+f) \circ(\operatorname{Id}+u+h)\right. \\
& \left.-D u^{n}(\omega+f) \circ(\operatorname{Id}+u)\right] \\
= & \mathbb{I}^{-}(I+D u)^{-1}[(\omega+f) \circ(\operatorname{Id}+u+h) \\
& -(\omega+f) \circ(\operatorname{Id}+u) \\
& \left.-D h(I+D u)^{-1}(\omega+f) \circ(\operatorname{Id}+u+h)\right] \\
& +\mathcal{O}\left(\|h\|_{\rho^{\prime}}^{\prime 2}\right) \\
= & \mathbb{I}^{-}(I+D u)^{-1}[D f \circ(\operatorname{Id}+u) h \\
& \left.-D h(I+D u)^{-1}(\omega+f) \circ(\operatorname{Id}+u)\right] \\
& +\mathcal{O}\left(\|h\|_{\rho^{\prime}}^{\prime 2}\right) .
\end{aligned}
$$

That completes the proof.
We want to find a solution for the equation

$$
\begin{equation*}
F(u)=0 \tag{2}
\end{equation*}
$$

For that, consider a continuous one-parameter family of maps: $U_{\lambda}=\operatorname{Id}+u_{\lambda}$, $\lambda \in[0,1]$, with "initial" condition $U_{0}=I d$, i.e. $u_{0}=0$, such that

$$
F\left(u_{\lambda}\right)=(1-\lambda) F\left(u_{0}\right)
$$

Differentiating the above equation in respect to $\lambda$, we get

$$
\begin{equation*}
D F\left(u_{\lambda}\right) \frac{d u_{\lambda}}{d \lambda}=-F(0) \tag{3}
\end{equation*}
$$

Remark 9.3 The derivative of $F$ at $u$ can be rewritten in the form

$$
D F(u) h=\mathbb{I}^{-}(D U)^{-1}\left[D X \circ U h-D h(D U)^{-1} X \circ U\right] .
$$

Evaluating on $v \circ U$, where $v$ is the vector field that generates $U$, i.e. $d U / d \lambda=$ $v \circ U$ and $v \circ U=D U v$, one gets

$$
\begin{aligned}
D F(u) v \circ U & =\mathbb{I}^{-}(D U)^{-1}[v, X] \circ U \\
& =\mathbb{I}^{-}\left[v,(D U)^{-1} X \circ U\right]
\end{aligned}
$$

where $[v, w]=D w v-D v w$ is the commutator for vector fields $v, w$. Thus, writing $\tilde{X}=(D U)^{-1} X \circ U$, it suffices to solve the commutator equation

$$
\mathbb{I}^{-}[v, \tilde{X}]=-\mathbb{I}^{-} X
$$

with respect to $v$ satisfying $\mathbb{I}^{+} v=0$, or, equivalently, invert $D F(u)$. Note also that we allow to have $v=v_{\lambda}$.

Proposition 9.4 If $u \in \mathcal{B}$, then $D F(u)^{-1}$ is a bounded linear operator from $\mathbb{I}^{-} \mathcal{A}_{d}\left(\rho^{\prime}\right)$ to $\mathbb{I}^{-} \mathcal{A}_{d}^{\prime}\left(\rho^{\prime}\right)$ and

$$
\left\|D F(u)^{-1}\right\|<\frac{\hat{\delta}}{\hat{\varepsilon}}
$$

From the above proposition (to be proved in Section 9.3) we know that $D F(u)$ is invertible for $u \in \mathcal{B}$, thus we may integrate (3) with respect to $\lambda$, obtaining the integral equation:

$$
\begin{equation*}
u_{\lambda}=-\int_{0}^{\lambda} D F\left(u_{\mu}\right)^{-1} F(0) d \mu \tag{4}
\end{equation*}
$$

In order to check that $u_{\lambda} \in \mathcal{B}$ for any $\lambda \in[0,1]$, we estimate its norm:

$$
\begin{aligned}
\left\|u_{\lambda}\right\|_{\rho^{\prime}}^{\prime} & \leq \sup _{v \in \mathcal{B}}\left\|D F(v)^{-1} F(0)\right\|_{\rho^{\prime}}^{\prime} \\
& \leq \sup _{v \in \mathcal{B}}\left\|D F(v)^{-1}\right\| \cdot\left\|\mathbb{I}^{-} f\right\|_{\rho^{\prime}} \\
& <\frac{\hat{\delta} \hat{\hat{\varepsilon}}}{\hat{\varepsilon}}=\hat{\delta}
\end{aligned}
$$

Therefore, the solution of (2) exists in $\mathcal{B}$ and is given by (4) when $\lambda=1$.
Now, the open ball $\hat{B}$ as claimed is simply given by $\hat{B}=\omega+\mathcal{E} \subset \mathcal{A}^{\prime}(\rho)$ using $\hat{\varepsilon}$ given by (6). For $X=\omega+f \in \hat{B}$ we have $\|u\|_{\rho^{\prime}}^{\prime}=\mathcal{O}\left(\|f\|_{\rho^{\prime}}\right)$. Thus,

$$
\begin{aligned}
\|\mathcal{U}(X)-\omega\|_{\rho^{\prime}} & \leq\left\|\sum_{n \geq 1}(-D u)^{n} \omega+(D U)^{-1} f \circ U\right\|_{\rho^{\prime}} \\
& \leq \frac{1}{1-\|u\|_{\rho^{\prime}}^{\prime}}\left(\|\omega\|\|u\|_{\rho^{\prime}}^{\prime}+\|f\|_{\rho}\right) \\
& \leq 2\left(\frac{\hat{\delta}}{\hat{\varepsilon}}\|\omega\|\|f\|_{\rho^{\prime}}+\|f\|_{\rho^{\prime}}\right) \\
& \leq 2\left(1+\max \left\{\frac{2(3-\sqrt{6})}{3}, \frac{6\|\omega\|}{\sigma}(\sqrt{6}+2)\right\}\right)\|X-\omega\|_{\rho}^{\prime}
\end{aligned}
$$

Notice that $\mathcal{U}(X)=(D U)^{-1}(\omega+f \circ U)$. Moreover, $\left\|\mathcal{U}(X)-\omega-\mathbb{I}^{+} f\right\|_{\rho^{\prime}}=$ $\mathcal{O}\left(\|f\|_{\rho^{\prime}}^{2}\right)$. This means that the derivative at $\omega$ of the map $\mathcal{U}$ is $\mathbb{I}^{+}$. Now, assume $\|f-\psi\|_{\rho}^{\prime}=\mathcal{O}\left(\|f\|_{\rho^{\prime}}\right), \omega+\psi \in \hat{B} \cap \mathbb{C}^{d}$, and write $\mathcal{U}(X)=\mathcal{U}(\omega+\psi+$ $(f-\psi))$. So,

$$
\left\|\mathcal{U}(\omega+\psi+(f-\psi))-(\omega+\psi)-\mathbb{I}^{+}(f-\psi)\right\|_{\rho^{\prime}}=\mathcal{O}\left(\|f\|_{\rho^{\prime}}^{2}\right)
$$

i.e. $\mathbb{I}^{+}$is the derivative of $\mathcal{U}$ at $\omega+\psi$. That completes the proof of Theorem 5.2.

### 9.3 Proof of Proposition 9.4

To prove Proposition 9.4 we start by inverting the derivative of $F$ at $u=0$.
Lemma 9.5 If $f \in \mathcal{E}$ such that $\|f\|_{\rho}^{\prime}<\sigma / 4$, the associated bounded linear operator $D F(0)^{-1}$, from $\mathbb{I}^{-} \mathcal{A}_{d}\left(\rho^{\prime}\right)$ to $\mathbb{I}^{-} \mathcal{A}_{d}^{\prime}\left(\rho^{\prime}\right)$, satisfies:

$$
\left\|D F(0)^{-1}\right\|<\frac{2}{\sigma-4\|f\|_{\rho}^{\prime}}
$$

Proof: From Lemma 9.2 one has

$$
\begin{aligned}
D F(0) h & =\mathbb{I}^{-}(\hat{f} h-D h \omega) \\
& =\mathbb{I}^{-}[\hat{f}-(D \cdot \omega)] h \\
& =-\left[\mathbb{I}-\mathbb{I}^{-} \hat{f}(D \cdot \omega)^{-1}\right](D \cdot \omega) h
\end{aligned}
$$

where $\hat{f} h=D f h-D h f$. Thus, the inverse of this operator, if it exists, is given by

$$
D F(0)^{-1}=-(D \cdot \omega)^{-1}\left[\mathbb{I}-\mathbb{I}^{-} \hat{f}(D \cdot \omega)^{-1}\right]^{-1}
$$

The inverse of $(D \cdot \omega)$ is the linear map from $\mathbb{I}^{-} \mathcal{A}_{d}\left(\rho^{\prime}\right)$ to $\mathbb{I}^{-} \mathcal{A}_{d}^{\prime}\left(\rho^{\prime}\right)$ :

$$
(D \cdot \omega)^{-1} g(\theta)=\sum_{k \in I^{-}} \frac{g_{k}}{2 \pi i k \cdot \omega} e^{2 \pi i k \cdot \theta}
$$

So,

$$
\begin{aligned}
\left\|(D \cdot \omega)^{-1} g\right\|_{\rho^{\prime}}^{\prime} & =\sum_{k \in I^{-}} \frac{1+2 \pi\|k\|}{2 \pi|k \cdot \omega|}\left\|g_{k}\right\| e^{\rho^{\prime}\|k\|} \\
& <\sum_{k \in I^{-}} \frac{1+2 \pi\|k\|}{2 \pi \sigma\|k\|}\left\|g_{k}\right\| e^{\rho^{\prime}\|k\|} \\
& \leq \frac{2}{\sigma}\|g\|_{\rho^{\prime}}
\end{aligned}
$$

where the use of the definition of $I^{-}$(Definition 5.1) was crucial to avoid dealing with arbitrarily small denominators. Hence, $\left\|(D \cdot \omega)^{-1}\right\|<\frac{2}{\sigma}$.

Similarly, for $\hat{f}: \mathbb{I}^{-} \mathcal{A}_{d}^{\prime}\left(\rho^{\prime}\right) \rightarrow \mathcal{A}_{d}\left(\rho^{\prime}\right)$,

$$
\|\hat{f} h\|_{\rho^{\prime}} \leq\|f\|_{\rho^{\prime}}^{\prime}\|h\|_{\rho^{\prime}}+\|f\|_{\rho^{\prime}}\|h\|_{\rho^{\prime}}^{\prime}
$$

Hence it is a bounded operator with $\|\hat{f}\| \leq 2\|f\|_{\rho^{\prime}}^{\prime}$. Therefore,

$$
\left\|\mathbb{I}^{-} \hat{f}(D \cdot \omega)^{-1}\right\|<\frac{4}{\sigma}\|f\|_{\rho^{\prime}}^{\prime}<1
$$

assuming $\|f\|_{\rho}^{\prime}<\sigma / 4$. So, $\left[\mathbb{I}-\mathbb{I}^{-} \hat{f}(D \cdot \omega)^{-1}\right]^{-1}$ is a bounded linear map in $\mathbb{I}^{-} \mathcal{A}_{d}\left(\rho^{\prime}\right)$ such that:

$$
\left\|\left[\mathbb{I}-\mathbb{I}^{-} \hat{f}(D \cdot \omega)^{-1}\right]^{-1}\right\|<\frac{\sigma}{\sigma-4\|f\|_{\rho^{\prime}}^{\prime}}
$$

completing the proof of the Lemma.
It remains now to estimate the variation of $D F$ with $u \in \mathcal{B}$.
Lemma 9.6 Let $0<\hat{\delta}<\min \left\{\frac{1}{2}, \frac{\rho-\rho^{\prime}}{4 \pi}\right\}$. Given $u \in \mathcal{B}$ such that $\|u\|_{\rho^{\prime}}^{\prime}=\delta<$ $\hat{\delta}$, the linear operator $D F(u)-D F(0)$, mapping $\mathbb{I}^{-} \mathcal{A}_{d}^{\prime}\left(\rho^{\prime}\right)$ into $\mathbb{I}^{-} \mathcal{A}_{d}\left(\rho^{\prime}\right)$ for $f \in \mathcal{E}$ with $\|f\|_{\rho}^{\prime}=\varepsilon<\hat{\varepsilon}$, is bounded and

$$
\|D F(u)-D F(0)\|<\frac{\delta}{1-\delta}\left[\left(\frac{4 \pi}{\rho-\rho^{\prime}}+\frac{4-2 \delta}{1-\delta}\right) \varepsilon+\frac{2-\delta}{1-\delta}\|\omega\|\right]
$$

Proof: Lemma 9.2 gives us that

$$
\begin{aligned}
{[D F(u)-D F(0)] \hbar=} & \mathbb{I}^{-}(I+D u)^{-1}[D f \circ(\operatorname{Id}+u) h-(I+D u) D f h \\
& -D h(I+D u)^{-1}(\omega+f) \circ(\operatorname{Id}+u) \\
& +(I+D u) D h(\omega+f)] \\
= & \mathbb{I}^{-}(I+D u)^{-1}\{[D f \circ(\operatorname{Id}+u)-D f-D u D f] h \\
& +D u D h(\omega+f) \\
& \left.-D h(I+D u)^{-1}[f \circ(\operatorname{Id}+u)-f-D u(\omega+f)]\right\} \\
= & \mathbb{I}^{-}(I+D u)^{-1}\{A+B+C\}
\end{aligned}
$$

where $A, B$ and $C$ are each of the respective three terms in the previous sum. Thus, $\left\|(I+D u)^{-1}\right\| \leq(1-\|D u\|)^{-1}$. Using Lemma 9.1, one obtains:

$$
\|A\|_{\rho^{\prime}} \leq\left(\frac{4 \pi}{\rho-\rho^{\prime}}\|f\|_{\rho}^{\prime}\|u\|_{\rho^{\prime}}+\|f\|_{\rho^{\prime}}^{\prime}\|u\|_{\rho^{\prime}}^{\prime}\right)\|h\|_{\rho^{\prime}}
$$

and $\|B\|_{\rho^{\prime}} \leq\left(\|\omega\|+\|f\|_{\rho^{\prime}}\right)\|u\|_{\rho^{\prime}}^{\prime}\|h\|_{\rho^{\prime}}^{\prime}$. Finally, noting that

$$
\left\|D u(I+D u)^{-1}\right\| \leq\|D u\| \cdot(1-\|D u\|)^{-1}
$$

and from Lemma 9.1 again, one gets the bound for the third term:

$$
\|C\|_{\rho^{\prime}} \leq \frac{1}{1-\|u\|_{\rho^{\prime}}^{\prime}}\left[\|f\|_{\left(\rho+\rho^{\prime}\right) / 2}^{\prime}\|u\|_{\rho^{\prime}}+\|u\|_{\rho^{\prime}}^{\prime}\left(\|\omega\|+\|f\|_{\rho^{\prime}}\right)\right]\|h\|_{\rho^{\prime}}^{\prime}
$$

To conclude the proof of Proposition 9.4, notice that, for $u \in \mathcal{B}$ and $f \in \mathcal{E}$ :

$$
\begin{aligned}
\left\|D F(u)^{-1}\right\| & \leq\left(\left\|D F(0)^{-1}\right\|^{-1}-\|D F(u)-D F(0)\|\right)^{-1} \\
& <\left\{\frac{\sigma}{2}-2 \hat{\varepsilon}-\frac{\hat{\delta}}{1-\hat{\delta}}\left[\left(\frac{4 \pi}{\rho-\rho^{\prime}}+\frac{4-2 \hat{\delta}}{1-\hat{\delta}}\right) \hat{\varepsilon}+\frac{2-\hat{\delta}}{1-\hat{\delta}}\|\omega\|\right]\right\}^{-1} \\
& <\frac{\hat{\delta}}{\hat{\varepsilon}} .
\end{aligned}
$$

The last inequality is true if we choose $\hat{\varepsilon}$ and $\hat{\delta}$ to satisfy

$$
\begin{equation*}
\hat{\varepsilon}<\hat{\delta}\left[\frac{\sigma}{2}-\frac{2 \hat{\delta}}{(1-\hat{\delta})^{2}}\|\omega\|\right]\left[1+2 \hat{\delta}+\frac{\hat{\delta}^{2}}{1-\hat{\delta}}\left(\frac{4 \pi}{\rho-\rho^{\prime}}+\frac{4-2 \hat{\delta}}{1-\hat{\delta}}\right)\right]^{-1} \tag{5}
\end{equation*}
$$

and

$$
\hat{\delta}<\frac{1}{2}-\frac{1}{2}\left[1+\frac{\sigma}{2\|\omega\|}\right]^{-\frac{1}{2}}<\frac{1}{2}
$$

An effective value of $\hat{\varepsilon}$ for each $\sigma$ can then be determined by the minimum value of the upper bounds given in Lemma 9.5 and (5). That can be obtained from a specific choice of $\hat{\delta}$ given by

$$
\hat{\delta}=\min \left\{\frac{\rho-\rho^{\prime}}{4 \pi}, \frac{3-\sqrt{6}}{6} \frac{\sigma}{\|\omega\|}\right\}<\frac{1}{2}-\frac{1}{2}\left[1+\frac{\sigma}{2\|\omega\|}\right]^{-\frac{1}{2}}<\frac{1}{2}
$$

as long as $\sigma /\|\omega\|<1$. Therefore, it is sufficient to impose

$$
\begin{equation*}
\hat{\varepsilon}:=\min \left\{\frac{\sigma}{4}, \frac{\sqrt{6}-2}{12} \sigma \hat{\delta}\right\} \tag{6}
\end{equation*}
$$

where the following inequalities were applied to (5):

$$
\frac{1}{1-\hat{\delta}}<\sqrt{1+\frac{\sigma}{2\|\omega\|}} \quad \text { and } \quad \frac{2 \hat{\delta}}{(1-\hat{\delta})^{2}}\|\omega\|<\frac{(3-\sqrt{6}) \sigma}{2}<\frac{\sigma}{2}
$$

### 9.4 Analyticity Improvement

We want to prove Proposition 6.4, i.e. that $X \circ T$ is analytic in $\mathcal{D}(\rho)$ and has bounded derivative. Let $\sigma>0$ such that

$$
\max \left\{\frac{\left\|T^{*} x\right\|}{\|x\|}:\|x\| \neq 0,|\omega \cdot x| \leq \sigma\|x\|\right\}<\kappa
$$

$\mathrm{Using} X \circ T(\theta)=\sum_{k \in I_{\sigma}^{+}} X_{k} e^{2 \pi i T^{*} k \cdot \theta}$,

$$
\begin{aligned}
\|X \circ T\|_{\rho} & \leq \sum_{k \in I_{\sigma}^{+}}\left\|X_{k}\right\| e^{\rho\left\|T^{*} k\right\|} \\
& \leq \sum_{k \in I_{\sigma}^{+}}\left\|X_{k}\right\| e^{\left(\rho \kappa-\rho^{\prime}\right)\|k\|} e^{\rho^{\prime}\|k\|} \leq\|X\|_{\rho^{\prime}}
\end{aligned}
$$

and for the derivative we have

$$
\begin{aligned}
\|D(X \circ T)\|_{\rho} & \leq 2 \pi \sum_{k \in I_{\sigma}^{+}}\left\|T^{*} k\right\| e^{-\beta\|k\|}\left\|X_{k}\right\| e^{\left(\rho \kappa+\beta-\rho^{\prime}\right)\|k\|} e^{\rho^{\prime}\|k\|} \\
& \leq \frac{2 \pi}{\beta} \kappa\|X\|_{\rho^{\prime}},
\end{aligned}
$$

by choosing $0<\beta<\rho^{\prime}-\kappa \rho$, and using the relation $\sup _{t>0} t e^{-\xi t} \leq 1 / \xi$ for $\xi>0$. These bounds imply that $\|\mathcal{T}(X)\|_{\rho}^{\prime} \leq(1+2 \pi \kappa / \beta)\left\|T^{-1}\right\|\|X\|_{\rho^{\prime}}$.

Let $r>\rho$ such that $\rho^{\prime}>r \kappa$. What was done above for $\mathcal{D}(\rho)$ applies as well to $\mathcal{D}(r)$. Therefore, one can decompose $\mathcal{T}=\mathcal{I} \circ \mathcal{J}$, where $\mathcal{J}: \mathbb{I}^{+} \mathcal{A}\left(\rho^{\prime}\right) \rightarrow$ $\mathcal{A}^{\prime}(r)$ as before, and $\mathcal{I}: \mathcal{A}^{\prime}(r) \rightarrow \mathcal{A}^{\prime}(\rho)$ is the inclusion map $\mathcal{I}(X)=\left.X\right|_{\mathcal{D}(\rho)}$. Note that $\mathcal{J}$ is bounded and $\mathcal{I}$ is compact, thus completing the proof.

## 10 Spectral Properties of $D \mathcal{R}(\omega)$ and Invariant Manifolds

As $\sigma$ remains fixed, we continue to drop the index of the sets $I^{ \pm}$and the projections $\mathbb{I}^{ \pm}$.

First we study the behaviour of the linear map $\mathcal{L}$ whose spectral properties are closely related to the ones of $D \mathcal{R}(\omega)$. Notice that $\mathbb{E} \mathcal{A}^{\prime}(\rho)=\mathbb{C}^{d}$ and $(\mathbb{I}-\mathbb{E}) \mathcal{A}^{\prime}(\rho)$ are invariant subspaces of $\mathcal{L}$ as $\mathcal{L} \circ \mathbb{E}=\mathbb{E} \circ \mathcal{L}$.

Lemma 10.1 If $\omega \in D C(\beta)$, for any $\beta>0$, and $\sigma\|\bar{\omega}\|<\frac{1}{2}$, one finds constants $a>0$ and $b, c>1$ such that

$$
\left\|\mathcal{L}^{n+1}(\mathbb{I}-\mathbb{E})\right\| \leq b \lambda_{1}^{2} e^{-a c^{n}}\left\|\mathcal{L}^{n}(\mathbb{I}-\mathbb{E})\right\|, \quad n \geq 0
$$

Proof: For every $n \geq 0$, we define $I_{n}^{+}$as the set of indices $k$ that verify $\left(T^{*}\right)^{n} k \in I^{+} \backslash\{0\}$ (see Figure 4). Hence,

$$
\sigma\left\|\left(T^{*}\right)^{n} k\right\| \geq\left|\omega \cdot\left(T^{*}\right)^{n} k\right|=\left|T^{n} \omega \cdot k\right|=\left|\lambda_{1}\right|^{n}|\omega \cdot k|, \quad \text { for } k \in I_{n}^{+}
$$

This inequality implies that

$$
|\omega \cdot k| \leq \sigma\|(\omega \cdot k) \bar{\omega}\|+\sigma\left\|\left(\lambda_{1}^{-1} T^{*}\right)^{n}[k-(\omega \cdot k) \bar{\omega}]\right\| .
$$

Thus, assuming $\sigma\|\bar{\omega}\|<\frac{1}{2}$ and knowing that $k-(\omega \cdot k) \bar{\omega}$ is the component of $k \in I_{n}^{+}$on the spectral directions of $T^{*}$ corresponding to the eigenvalues $\lambda_{j}, j=2, \ldots, d$,

$$
|\omega \cdot k| \leq \frac{\sigma}{1-\sigma\|\bar{\omega}\|}\left|\frac{\lambda_{2}}{\lambda_{1}}\right|^{n}\|k-(\omega \cdot k) \bar{\omega}\|
$$

As $\omega$ is a diophantine vector of degree $\beta>0$, there is a constant $C>0$ such that

$$
|\omega \cdot k|>C\|k\|^{-\beta-1}
$$



Figure 4: Set $I_{n}^{+} \subset I^{+}$for the two-dimensional case.
and we get a lower bound for the norm of $k \in I_{n}^{+}$,

$$
\begin{equation*}
\|k\|^{2+\beta}>C \frac{1-\sigma\|\bar{\omega}\|}{2 \sigma}\left|\frac{\lambda_{1}}{\lambda_{2}}\right|^{n} \tag{7}
\end{equation*}
$$

Let us now define the operator $\mathbb{I}_{n}^{+}: \mathcal{A}^{\prime}(r) \rightarrow \mathcal{A}^{\prime}(\rho)$, where we have chosen $r>\rho$. This is a projection for the indices in $I_{n}^{+}$together with an analytic inclusion. Making use of the bound (7), the operator norm of $\mathbb{I}_{n}^{+}$follows from

$$
\left\|\mathbb{I}_{n}^{+} f\right\|_{\rho}^{\prime}=\sum_{I_{n}^{+}}(1+2 \pi\|k\|)\left\|f_{k}\right\| e^{r\|k\|} e^{-(r-\rho)\|k\|} \leq e^{-a c^{n}}\|f\|_{r}^{\prime}
$$

with $a=(r-\rho)\left[C(2 \sigma)^{-1}(1-\sigma\|\bar{\omega}\|)\right]^{1 /(2+\beta)}>0$ and $c=\left|\lambda_{1} \lambda_{2}^{-1}\right|^{1 /(2+\beta)}>1$.
We can write $\mathcal{L}^{n+1}$ with respect to $\mathcal{L}^{n}$, including the operator $\mathbb{I}_{n}^{+}$, in the form

$$
\mathcal{L}^{n+1}(\mathbb{I}-\mathbb{E}) f=\mathcal{L}^{n}(\mathbb{I}-\mathbb{E}) \mathbb{I}_{n}^{+} \overline{\mathcal{L}} f
$$

where $\overline{\mathcal{L}}: \mathcal{A}^{\prime}(\rho) \rightarrow \mathcal{A}^{\prime}(r)$ is $\mathcal{L}$ followed by an analytic extension. That is, $\overline{\mathcal{L}}=\lambda_{1} \overline{\mathcal{I}} \circ \mathbb{I}^{+}$where $\overline{\mathcal{T}}: \mathbb{I}^{+} \mathcal{A}^{\prime}(\rho) \subset \mathbb{I}^{+} \mathcal{A}(\rho) \rightarrow \mathcal{A}^{\prime}(r)$ given as in Proposition 6.4. The norm $\|\overline{\mathcal{L}}\| \leq\|\overline{\mathcal{T}}\|$ is determined by

$$
\|\overline{\mathcal{T}} f\|_{r}^{\prime} \leq b\left\|T^{-1}\right\|\|f\|_{\rho}^{\prime}
$$

for some $r>\rho$ and a constant $b>1$ from the proof of Proposition 6.4 in Section 9.4. So, the claim follows from

$$
\left\|\mathcal{L}^{n+1}(\mathbb{I}-\mathbb{E}) f\right\|_{\rho}^{\prime} \leq\left\|\mathcal{L}^{n}(\mathbb{I}-\mathbb{E})\right\|\left\|\mathbb{I}_{n}^{+}\right\|\|\overline{\mathcal{L}}\|\|f\|_{\rho}^{\prime}
$$

and $\left\|T^{-1}\right\| \leq\left|\lambda_{1}\right|$.
So, the direction given by the non-constant terms $\mathbb{I}-\mathbb{E}$ is a stable eigenspace with eigenvalue zero. The spectrum of $\mathcal{L} \circ \mathbb{E}$ is simply the one of $T^{-1}$. That is, the eigenvalues are $\lambda_{1} \lambda_{j}^{-1}, j=2, \ldots, d$, and 1 corresponding to the eigenvector $\omega$. Hence, there is a "neutral" direction given by $\omega$ whereas the remaining ones corresponding to $\omega^{(2)}, \ldots, \omega^{(d)}$ are all unstable.

The spectral properties of $D \mathcal{R}(\omega)$ are easily related to those of $\mathcal{L}$ because

$$
D \mathcal{R}(\omega)=\left(\mathbb{I}-\mathbb{P}_{\omega} \mathbb{E}\right) \mathcal{L}
$$

where $\mathbb{P}_{\omega}$ is the projection on the subspace spanned by $\omega$, i.e. $\mathbb{P}_{\omega} f_{0}=$ $\left(\omega \cdot f_{0}\right) \bar{\omega}$. Note also that $\mathcal{L}$ commutes with the projection $\mathbb{P}_{\omega} \mathbb{E}$ because $\omega$ is a eigenvector of $T$, thus of $\mathcal{L}$. So,

$$
\left(\mathbb{I}-\mathbb{P}_{\omega} \mathbb{E}\right) \mathcal{L}=\mathcal{L}\left(\mathbb{I}-\mathbb{P}_{\omega} \mathbb{E}\right)
$$

Therefore, the projection $\mathbb{P}^{u}=\left(\mathbb{I}-\mathbb{P}_{\omega}\right) \mathbb{E}$ determines the linear space $\mathbb{P}^{u} \mathcal{A}^{\prime}(\rho)$ spanned by $\left\{\omega^{(2)}, \ldots \omega^{(j)}\right\}$, which is the unstable eigenspace of the linearised map at $\omega$. In a similar way, $\mathbb{P}^{s}=\mathbb{I}-\mathbb{P}^{u}$ is the stable part.

We are now in a position to define and determine the stable and unstable invariant local manifolds of $\mathcal{R}$ at the fixed point $\omega$. The stable one, $\mathcal{W}^{s}(\omega)$, corresponds to all the vector fields $X$ in some neighbourhood $B$ such that $\mathcal{R}^{n}(X) \in B$, for any $n \geq 0$. Similarly, considering the inverse iterates of $\mathcal{R}$, we define the unstable local manifold $\mathcal{W}^{u}(\omega)$.

If $X \in \mathcal{W}^{s}(\omega)$, the winding set is equal to that of $\omega$, since it is preserved by $\mathcal{R}$. This is in fact true because $\omega$ is an eigenvector of the matrix $T \in$ $G L(d, \mathbb{Z})$, the linear part of the transformation induced by $\mathcal{R}$ on the domain of $X$. The winding set is preserved up to the action of $\hat{T}: x \mapsto T x /\|T x\|$, by iterating $\mathcal{R}$. If the winding set of a vector $X \in B$ is $w_{X}=\omega /\|\omega\|$, then it is also in $\mathcal{W}^{s}(\omega)$, otherwise would approximate the unstable direction and $w_{X}$ would be different.

The set of constant vector fields is invariant under $\mathcal{R}$, i.e. $\mathcal{R}\left(\mathbb{C}^{d} \cap B\right) \subset$ $\mathbb{C}^{d}$, as can be seen from the way the non-linear coordinate change $\mathcal{U}$ was constructed in Theorem 5.2. This allows us to conclude that the invariant unstable local manifold is $\mathcal{W}^{u}(\omega)=\mathbb{P}^{u}(B)$.

Lemma 10.2 We can find a ball $B^{\prime} \subset \mathcal{A}_{d}^{\prime}(\rho)$ around $\omega$ such that, for any $X \in \mathcal{W}^{s}(\omega) \cap B^{\prime}$, there are constants $K>0$ and $\left.\theta \in\right] 0,1[$ yielding

$$
\left\|\mathcal{R}^{m}(X)-\omega\right\|_{\rho}^{\prime}<K \theta^{m}\|X-\omega\|_{\rho}^{\prime}, \quad m>0
$$

Moreover, there exists $a>0, c>1$ and $N>0$ satisfying

$$
\left\|(\mathbb{I}-\mathbb{E}) \mathcal{R}^{n}(X)\right\|_{\rho}^{\prime}<e^{-a c^{n}}\|X-\omega\|_{\rho}^{\prime}+\theta^{2 n}\|X-\omega\|_{\rho}^{2}, \quad n>N
$$

Proof: Fixing $0<t^{\prime}<t<\theta<1$ and choosing an appropriate norm $\|\cdot\|$ on $\mathcal{A}^{\prime}(\rho)$ equivalent to $\|\cdot\|_{\rho}^{\prime}$, it is possible to have $\|D \mathcal{R}(\omega) f\|<t^{\prime}\|f\|$, $f \in \mathbb{P}^{s}\left(\mathcal{A}^{\prime}(\rho)\right)$, since $D \mathcal{R}(\omega)$ is compact.

We can write $X=\omega+f+g$ where $f=\mathbb{P}^{s}(X-\omega)$ and $g=\mathbb{P}^{u}(X-\omega)$. By the stable manifold theorem, there is a constant $A>0$ for which

$$
\begin{equation*}
\|g\|<A\|f+g\|^{2}, \quad \text { and } \quad\|f\|<\|f+g\|+A\|f+g\|^{2} \tag{8}
\end{equation*}
$$

The image of $X$ is then $X^{\prime}=\mathcal{R}(X)=\omega+f^{\prime}+g^{\prime}$ with

$$
\begin{align*}
f^{\prime} & =D \mathcal{R}(\omega) f \\
g^{\prime} & =D \mathcal{R}(\omega) g+\mathcal{O}\left(\|f+g\|^{2}\right) \tag{9}
\end{align*}
$$

The analyticity of $\mathcal{R}$ and the formula of the second-order Taylor remaining imply that there exists constants $B, C>0$ satisfying

$$
\begin{align*}
& \left\|f^{\prime}\right\|<t^{\prime}\|f\| \\
& \left\|g^{\prime}\right\|<C\|g\|+B\|f+g\|^{2} \tag{10}
\end{align*}
$$

By combining this with (8) we obtain

$$
\left\|f^{\prime}+g^{\prime}\right\|<t^{\prime}\|f+g\|+\left(A t^{\prime}+A C+B\right)\|f+g\|^{2}
$$

So, there is a radius $r>0$ such that

$$
\left\|X^{\prime}-\omega\right\|<t\|X-\omega\|
$$

for $\|X-\omega\|<r$. Iterating this inequality $m$ times we obtain

$$
\left\|\mathcal{R}^{m}(X)-\omega\right\|_{\rho}^{\prime}<K \theta^{m}\|X-\omega\|_{\rho}^{\prime}
$$

if $B^{\prime}=\left\{X \in \mathcal{A}_{d}^{\prime}(\rho):\|X-\omega\|<r\right\}$, that is the first of our claims.
From Lemma 10.1 we can find $a>0$ and $c>1$ and a sequence of positive integers $k_{n}<\varepsilon n, \varepsilon>0$, such that

$$
\left\|(\mathbb{I}-\mathbb{E})[D \mathcal{R}(\omega)]^{k_{n}}\right\|<e^{-a c^{n}}
$$

whenever $n$ is sufficiently large. Given such $n$ and $X \in \mathcal{W}^{s}(\omega) \cap B^{\prime}$, the first part of the proof showed that $\left\|\mathcal{R}^{m}(X)-\omega\right\|<\delta=t^{m} r$, with $m=n-k_{n}$. Consider a map $F:(f, g) \mapsto\left(f^{\prime}, g^{\prime}\right)$ defined by $(9)$, and let $f_{0}=\mathbb{P}^{s}\left(\mathcal{R}^{m}(X)-\right.$ $\omega)$ and $g_{0}=\mathbb{P}^{u}\left(\mathcal{R}^{m}(X)-\omega\right)$. The pair $\left(f_{0}, g_{0}\right)$ satisfies

$$
\begin{aligned}
& \left\|f_{0}\right\|<2 \delta \\
& \left\|g_{0}\right\|<A \delta^{2}
\end{aligned}
$$

Applying the inequalities (10), if $1<A<1 / \delta$, then $\left\|g_{1}\right\|<(C+9 B) A \delta^{2}$, using the notation $\left(f_{i}, g_{i}\right)=F^{i}\left(f_{0}, g_{0}\right), i \geq 0$. The fact that $\left\|f_{i}\right\|$ is bounded by at least $2 \delta$ allows us to iterate $F$ for $k_{n}$ times, obtaining:

$$
\begin{aligned}
& f_{k_{n}}=[D \mathcal{R}(\omega)]^{k_{n}} f_{0} \\
& \left\|g_{k_{n}}\right\|<(C+9 B)^{k_{n}} A \delta^{2}
\end{aligned}
$$

as long as $(C+9 B)^{k_{n}} A \delta<1$. This is in fact verified if we choose $\varepsilon>0$ such that $(C+9 B)^{k_{n}} t^{2 m}<\theta^{2 n}$, or, more strongly,

$$
\left[t^{-2}(C+9 B)\right]^{\varepsilon}<(\theta / t)^{2}
$$

Now, $\mathcal{R}^{n}(X)=\omega+f_{k_{n}}+g_{k_{n}}$ and $(\mathbb{I}-\mathbb{E}) \mathcal{R}^{n}(X)=(\mathbb{I}-\mathbb{E})\left(f_{k_{n}}+g_{k_{n}}\right)$. Hence,

$$
\left\|(\mathbb{I}-\mathbb{E}) \mathcal{R}^{n}(X)\right\|<2 e^{-a c^{n}} t^{m} r+A \theta^{2 n} r^{2}
$$

which completes the proof, if $n$ is chosen sufficiently large.

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