Gevrey estimates for asymptotic expansions of Tori in weakly dissipative systems

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Received 15 January 2021, revised 30 December 2021
Accepted for publication 11 March 2022
Published 25 May 2022

Abstract
We consider a singular perturbation for a family of analytic symplectic maps of the annulus possessing a KAM torus. The perturbation introduces dissipation and contains an adjustable parameter. By choosing the adjustable parameter, one can ensure that the torus persists under perturbation. Such models are common in celestial mechanics. In field theory, the adjustable parameter is called the counterterm and in celestial mechanics, the drift. It is known that there are formal expansions in powers of the perturbation both for the quasi-periodic solution and the counterterm. We prove that the asymptotic expansions for the quasiperiodic solutions and the counterterm satisfy Gevrey estimates. That is, the $n$th term of the expansion is bounded by a power of $n!$. The Gevrey class (the power of $n!$) depends only on the Diophantine condition of the frequency and the order of the friction coefficient in powers of the perturbative parameter. The method of proof we introduce may be of interest beyond the problem considered here. We consider a modified Newton method in a space of power expansions. As is customary in KAM theory, each step of the method is estimated in a smaller domain. In contrast with the KAM results, the domains where we control the Newton method shrink very fast and the Newton method does not prove that the solutions are analytic. On the other hand, by examining carefully the process, we can obtain estimates on the coefficients of the expansions and conclude the series are Gevrey.

Keywords: Gevrey estimates, dissipative systems, quasi-periodic solutions
Mathematics Subject Classification numbers: 35C20, 70K70, 70K43, 37J40, 34K26, 30E10.

Recommended by Dr Richard Montgomery.
* Both authors have been supported by NSF Grant DMS 1800241.
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1. Introduction

Hamiltonian systems with small dissipation appear as models of many problems of physical interest. Notably, dissipation is a small effect in astrodynamics of planets and satellites [MNF87, Cel13]. In the design of many mechanical devices, eliminating friction is a design goal which is never completely accomplished. Hamiltonian systems with friction also appear as Euler–Lagrange equations of discounted functionals which are natural in finance and in the receding horizon problem in control theory. In such a case the limit of zero discount (equivalent to the limit of zero friction) is of interest. See [Ben88, DFIZ16, ISM11, MHER95] for different studies of the zero dissipation limit in calculus of variations and in control.

When the friction is small, it is natural to study such systems using perturbation theory. Nevertheless, adding a small friction is a very singular perturbation, and periodic/quasi-periodic orbits may disappear for arbitrarily small values or the perturbation. In contrast with Hamiltonian systems that often have sets of quasi-periodic orbits of positive measure (KAM theorem), for dissipative forced systems, there are few periodic or quasi-periodic orbits. These quasi-periodic orbits of a fixed frequency are known to persist only if one can adjust parameters in the system [BHS96, Mos67, Sev99]. As discussed very clearly in [Mos73], the number of parameters needed is affected by the geometric properties of the systems considered.

In recent times, for some particular types of dissipative systems—the conformally symplectic systems, see definition 1—there is a very systematic KAM theory [CCdlL13] based on geometric arguments. The examples mentioned above (Hamiltonian systems with friction proportional to the momentum and Euler–Lagrange equations of exponentially discounted variational principles) are conformally symplectic. This theory, once we fix a frequency, predicts the changes of parameters and the changes in the solutions needed to obtain a quasi-periodic solution of the prescribed frequency.

The goal of this paper is to study the singular perturbation theories in which the small parameter also introduces dissipation.

There are several studies of the singular perturbation theories in dissipation which are particularly relevant for us: the paper [CCdlL17] shows that if one fixes a Diophantine frequency \( \omega \) (see definition 11), considers a Hamiltonian system—not necessarily integrable—with a quasi-periodic solution of frequency \( \omega \), and introduces a conformally symplectic perturbation (see definition 1), then there is a (unique under a natural normalization) formal power series expansion for the quasi-periodic solution of frequency \( \omega \) and for the drift parameter. These series are very similar to the Lindstedt series of classical mechanics. The paper [CCdlL17] also showed that the formal Lindstedt series is the asymptotic expansion of a true solution defined in a complex domain of perturbation parameters that does not include any ball around zero (giving an indication that the power series may be divergent). The paper [BC19] studied numerically these Lindstedt series in a concrete example and the possible domain of analyticity of the function (using Padé as well as non-perturbative methods). The numerical studies in [BC19] lead to the remarkable conjecture that, in the cases examined, the formal power series giving the quasiperiodic solution and the forcing are Gevrey (see definition 9).

In this paper, for some class of analytic maps (we require that the system is conformally symplectic and that the nonlinearity is a trig. polynomial) we show that the conjecture in [BC19] is true and that the series obtained are indeed Gevrey. The Gevrey class can be bounded.

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1 A problem in astrodynamics which motivate us is the spin orbit problem describing approximately the motion of an oblate planet, subject to tidal friction, in a Keplerian orbit [Cel91].
depending only on the Diophantine condition of the frequency $\omega$ (and the order of the friction in the dissipation). See theorem A.

The Gevrey class of functions has received a lot of interest recently since those functions are related to many deep theorems of dynamical systems (KAM, Nekhoroshev). Similar theories (e.g. hypoellipticity) also admit Gevrey classes as natural regularity. This paper goes in a different direction. Even if we start with an analytic problem—indeed polynomial!—several objects of interest are only Gevrey. The phenomenon that analytic problems have only Gevrey solutions has appeared in other contexts in dynamics, notably in the study of singular perturbations [CDRSS00], the regularity of attractors and fast-slow systems [Bae95, CD91, FT89].

Closer to us, in dependence on parameters of solutions of nonlinear problems, [Lin92, Sau92], dependence of KAM tori in the frequency [Pop00], or in the theory of parabolic manifolds [BFM17, BH08].

We note that showing that a perturbative expansion is Gevrey allows to obtain good bounds of the error of a finite sum [BDM08]. It also allows the use of resummation methods to extract better results for the series, [CGGG07], and it gives insights on the analyticity domains. Indeed, in the mathematical physics literature, there has been considerable interest in the Gevrey nature of perturbation theories, often called factorial bounds, Borel summability, etc [FMRS87, GBD05, GG84]. We hope that introducing a new method for these questions can have interest in other motivations.

The method of proof we introduce may be of interest beyond the problem considered here and we hope that there are other applications. We consider a Newton method in the space of power expansions. As in KAM theory, each step of the quadratically convergent method is estimated in a domain smaller than the domain of the previous steps. In contrast with KAM theory, the domains where we control the results shrink very fast to a point, so that, at the end we do not obtain any analytic function. On the other hand, by examining carefully the process, we obtain estimates on the coefficients of the expansions.

Our hypothesis that the nonlinearity is a trigonometric polynomial ensures that the coefficients of order $N$ do not change after $\log_2(N)$ steps of the Newton method, so that one can use Cauchy estimates in the domain that is under control after $\log_2(N)$ steps to obtain estimates on the $N$th coefficient.

We hope that the hypothesis that the nonlinearity is a trigonometric polynomial can be removed at the price of estimating the change of the coefficients in subsequent iterations, but a proof would require a new set of estimates that—if indeed possible—would lengthen the exposition and obscure the main ideas.

The Newton method acting on power series is patterned after the Newton method used in [CCdlL13]. This Newton method takes advantage of remarkable cancellations related to the geometry and introduces the corrections to the torus additively (rather than making changes of variables). The fact that the Newton method in [CCdlL13] does not involve changes of variables makes it possible to lift it to formal power series. We will present full details later.

For simplicity in the treatment, we will deal with maps since the geometric arguments are simpler. The same arguments apply for differential equations, but they are more elaborate. Besides adapting the proof of maps to the case of ODE’s, one can deduce rigorously the results for differential equations from the results for maps by taking time-$T$ maps. Note that in this case, the fact that the nonlinearity in the time-$T$ map is a trig. polynomial is difficult to express in terms of the original ODE. This is another reason why we would like eventually to get rid of that hypothesis.
1.1. A preview of the main results

A model to keep in mind is the so-called dissipative standard map $f_{\varepsilon, \mu_e} : T \times \mathbb{R} \to T \times \mathbb{R}$ given by

$$f_{\varepsilon, \mu_e}(x, y) = \left( x + \lambda(\varepsilon)y + \mu_e - \varepsilon V(x), \lambda(\varepsilon)y + \mu_e - \varepsilon V(x) \right).$$  \hfill (1.1)

In (1.1), the physical meaning of $\lambda(\varepsilon) = 1 - \varepsilon^\alpha$, $\alpha \in \mathbb{N}$, is dissipation and $\mu_e$, called the drift parameter, has the physical meaning of a forcing. Our assumption on the nonlinearity amounts to $V$ being a trigonometric polynomial. The model (1.1) is indeed conformally symplectic in the sense of definition 1 (see below). The map (1.1) is the model that was used in the numerical experiments in [BC19]. The model (1.1) can be thought as a numerical time step—using a Verlet-like method—of the spin–orbit problem

$$\begin{align*}
\dot{x} &= y \\
\dot{y} &= -\mu y + \lambda + \varepsilon V(x).
\end{align*}$$  \hfill (1.2)

Note that for $\varepsilon = 0$ and $\mu_0 = 0$, the map (1.1) is integrable. The integrability of the map at $\varepsilon = 0$ does not play any role in the theoretical results in [CCdlL17], the only assumption needed in [CCdlL17] is that the map for $\varepsilon = 0$ is symplectic and has an invariant torus. For the numerical study in [BC19], the fact that the map for $\varepsilon = 0$ is integrable leads to much more efficient algorithms. In this paper, we will not use explicitly the integrability for $\varepsilon = 0$, but this seems to be the only case where it is possible to verify the assumption on the nonlinearity being a trig polynomial (yet another reason to try to get rid of that hypothesis).

The main results of this paper are theorems A–C. Theorems A and B establish the Gevrey character of the formal expansions for the drift parameter $\mu_e$ and for the quasi-periodic orbit of fixed frequency $\omega$, $K_{\varepsilon, \omega}$ of the map (1.1). The rigorous formulation of these theorems is given in section 3, the statements of the main results can be better understood after some preliminary definitions and remarks are given (see section 2). Here we give an informal statement summarizing our main results:

Given a Diophantine frequency $\omega$, the coefficients of the formal power series expansions $\sum K_n \varepsilon^n$ and $\sum \mu_n \varepsilon^n$ of the quasi-periodic orbit and the drift parameter, respectively, satisfy the following Gevrey estimates

$$\|K_n\| \leq C R^n n^{\tau/\alpha}$$  \hfill (1.3)

and

$$|\mu_n| \leq C R^n n^{\tau/\alpha},$$

where $\tau$ depends on the Diophantine type of $\omega$ (see definition 1) and $\alpha$ is the order of the dissipation $\lambda(\varepsilon) = 1 - \varepsilon^\alpha$. Moreover, the formal series $\sum K \varepsilon^j$ and $\sum \mu \varepsilon^j$ are asymptotic expansions of geometric objects $K_{\varepsilon, \omega}$ and $\mu_{\varepsilon, \omega}$ with respect to $\varepsilon$. Quantitative estimates are formulated in theorem C.

1.2. Organization of the paper

The paper is organized as follows. In section 2 we collect some standard definitions and we also define the function spaces in which the iterative procedure takes place. Also, in the same section we present some geometric identities which allow us to solve the linearized equations of the modified Newton method. In section 3 we state theorems A–C, which are the main results of the paper and establish the Gevrey character of the formal expansions of the quasi-periodic orbits.
The proof of theorem A is a consequence of theorem B. The proof of theorem B is done through the use of a quasi Newton method. In section 4 we formulate the iterative step of this Newton method, while in section 5 we provide estimates for the corrections and the new error at one step of the method. In section 6, using a KAM like argument, we give estimates for any step of the Newton like procedure and, with them, the proofs of theorems B and C are given. Finally, in the appendix A, we verify that the hypothesis HTP1 and HTP2 of theorem B are satisfied by the family of maps considered in theorem A.

2. Preliminaries

In this section we introduce the notations, collect some standard definitions including the Banach spaces and their norms that enter in this paper. This section should be used as a reference.

2.1. Symplectic properties

Let $\mathcal{M} = \mathbb{T}^d \times B, B \subseteq \mathbb{R}^d$, endowed with an exact symplectic form $\Omega$. Note that the manifold $\mathcal{M}$ is Euclidean (i.e. the tangent bundle is trivial) and we can compare vectors in different tangent spaces. This is crucial in KAM theory.

We denote by $J$ the matrix associated to the symplectic form $\Omega$, i.e. in coordinates we have $\Omega_x(u, v) = (u, J(x)v)$ where $(\cdot, \cdot)$ denotes the inner product for any $u, v \in T_x \mathcal{M}$. Note that $J$ depends on the choice of the inner product.

**Definition 1.** We say that a diffeomorphism defined on a symplectic manifold $(\mathcal{M}, \Omega)$ is conformally symplectic when $f^* \Omega = \lambda \Omega$ for a number $\lambda$, where $f^*$ denotes the standard pull back on forms.

The map (1.1) is conformally symplectic with the conformal factor $\lambda(\epsilon) = 1 - \epsilon^\alpha$ and the standard symplectic form $\Omega = dx \wedge dy$ on the cylinder $\mathbb{T} \times \mathbb{R}$.

2.2. Banach spaces of analytic functions

2.2.1. Analytic functions on the torus. Given $\rho > 0$ we define the complex extension of the $d$-dimensional torus as

$$
\mathbb{T}_\rho^d = \{ z \in \mathbb{C}^d / \mathbb{Z}^d \mid \text{Re}(z_j) \in \mathbb{T}, \text{Im}(z_j) \leq \rho \}
$$

and denote by $A_\rho$ as the $\mathbb{C}$-vector space of analytic functions defined in $\mathbb{T}_\rho^d$ which can be extended continuously to the boundary of $\mathbb{T}_\rho^d$. $A_\rho$ is endowed with the norm

$$
\|g\|_\rho = \sup_{\theta \in \mathbb{T}_\rho^d} \|g(\theta)\|
$$

which makes it into a Banach space.

For vector valued functions, $g = (g_1, g_2, \ldots, g_d)$, we define the norm

$$
\|g\|_\rho = \sqrt{\|g_1\|_\rho^2 + \|g_2\|_\rho^2 + \cdots + \|g_d\|_\rho^2}
$$
and for $d_1 \times d_2$ matrix valued functions, $G$, we define

$$
\|G\|_\rho = \sup_{v \in \mathbb{R}^{d_2}, \|v\|=1} \left( \sum_{i=1}^{d_1} \left( \sum_{j=1}^{d_2} |G_{ij}|^\rho v_j \right)^2 \right)^{1/2}.
$$

We will also need to work with functions of two variables. Denoting $B_\gamma(0) \subseteq \mathbb{C}$ the open ball with center zero and radius $\gamma$ in the complex plane, define

$$
A_{\rho, \gamma} = \{ K : B_\gamma(0) \to A_\rho \mid K \text{ is analytic in } B_\gamma(0) \times \text{ and can be extended continuously to } B_\gamma(0) \}
$$

endowed with the norm

$$
\|K\|_{\rho, \gamma} = \sup_{|\varepsilon| \leq \gamma} \|K(\varepsilon)\|_\rho.
$$

It is well known that with the norms $\|\cdot\|_{\rho, \gamma}$ and $\|\cdot\|_\rho$ the spaces $A_{\rho, \gamma}$ and $A_\rho$ are Banach algebras.

To discuss analyticity properties, we will need to deal with complex values of all the arguments. For physical applications, we need mainly real variables. Hence, it will be important that the functions we consider have the property that they yield real values for real arguments. The functions that satisfy this property (real valued for real arguments) is a closed (real) subspace of the above Banach spaces. All the constructions we use have the property that when applied to real valued functions, they produce real valued functions.

Note that we can think of functions $A_{\rho, \gamma}$ as analytic functions on $B_\gamma(0)$ taking values on a space of analytic functions of the torus. This point of view is consistent with the interpretation that we are considering families of problems and we are seeking families of solutions.

For typographical reasons from now on we will use the following notation. Given $K \in A_{\rho, \gamma}$ we use the notation $K(\epsilon)(\theta) = K(\epsilon, \theta) := (K(\varepsilon))(\theta)$.

**Definition 2.** Let $B$ a Banach space. Given an analytic function $g : B_\gamma(0) \subseteq \mathbb{C} \to B$, and $n \geq 0$, we say $g(\varepsilon) \sim O(\|\varepsilon\|^n)$ if and only if there exists $C > 0$ such that

$$
g(\varepsilon) \lesssim C|\varepsilon|^n
$$

for $\varepsilon$ small enough. Equivalently, $g(\varepsilon) \sim O(\|\varepsilon\|^n)$ if and only if $g(\varepsilon) = \sum_{k=0}^{\infty} g_k \varepsilon^k$ for $\varepsilon$ small enough and $g_k \in B$.

### 2.2.2. Cauchy estimates.

We recall the classical Cauchy inequalities, see [SZ65].

**Lemma 3.** Given any $0 < \delta \leq \rho$ and any function $f \in A_\rho$ with range in $\mathbb{C}^2 d$, then there exists a constant $C = C(d)$ such that

$$
\|Df\|_{\rho-\delta} \lesssim C\delta^{-1} \|f\|_\rho,
$$

where $Df$ denotes the derivative of $f$. We also have

$$
|\tilde{f}_k| \lesssim e^{-2\varepsilon k|\rho|} \|f\|_\rho,
$$

where $\tilde{f}_k$ are the components of $f$. 

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where $|k| = |k_1| + |k_2| + \cdots + |k_n|$ and $\hat{f}_k$ denotes the Fourier coefficient of $f$ with index $k$.

Remark 4. We point out that inequality (2.1), comes from the classical Cauchy inequality

$$\|\partial_\theta g\|_{\rho-\delta} \leq \delta^{-1} \|g\|_{\rho}$$

for any scalar function $g \in \mathcal{A}_\rho$. If $f$ is a vector valued function, as in lemma 3, the inequality (2.1) follows from the definitions of the norms at the beginning of the section.

As mentioned above we will be working with functions depending upon two variables. The following are Cauchy inequalities in the second variable, $\varepsilon$.

Lemma 5. For any $0 < r \leq \gamma$ and any function $f \in \mathcal{A}_{\rho,\gamma}$ such that $f_\varepsilon(\theta) = \sum_{n=0}^{\infty} f_n(\theta)\varepsilon^n$ we have

$$\|f_n\|_{\rho} \leq \frac{1}{n!} \|f\|_{\rho,r}.$$  

Proof. By Cauchy integral formula

$$f_n(\theta) = \frac{1}{n!} \frac{d^n}{d\varepsilon^n} f(\theta, \varepsilon) \bigg|_{\varepsilon=0} = \frac{1}{2\pi i} \int_{|\xi|=r} \frac{f(\theta, \xi)}{\xi^{n+1}} d\xi = \frac{1}{2\pi r^n} \int_0^{2\pi} f(\theta, re^{i\phi}) \frac{e^{in\phi}}{e^{i\varepsilon\phi}} d\phi,$$

thus, $|f_n(\theta)| \leq \frac{1}{n!} \sup_{|\varepsilon| \leq r} |f(\theta, \varepsilon)|$ and $\|f_n\|_{\rho} \leq \frac{1}{n!} \|f\|_{\rho,r}$. □

Corollary 6. Assume that $\Delta \in \mathcal{A}_{\rho,\gamma}$ is such that $\Delta_\varepsilon = \sum_{n=N+1}^{\infty} \Delta_n \varepsilon^n$. Let $a, b \in \mathbb{N}$ such that $N \leq a \leq b \leq \infty$ and denote $\Delta_\varepsilon^{(a,b)} = \sum_{n=a+1}^{b} \Delta_n \varepsilon^n$. Then, for all $0 < r < 1$ we have

$$\|\Delta_\varepsilon^{(a,b)}\|_{\rho,r} \leq \frac{r^{b+1}}{1-r} \|\Delta\|_{\rho,\gamma}.$$  

Remark 7. Note that the estimate in corollary 6 only depends on $a$, associated with the order of the first term in the expansion of $\Delta^{(a,b)}$.

2.3. Formal power series

2.3.1. General definitions. Formal power series expansions are just expressions of the form

$$\sum_n a_n \varepsilon^n,$$

where $a_n$ belong to a Banach space, sometimes $a_n$ are just scalars.

Formal power series are not meant to converge nor to represent a function. However, when the coefficients $a_n$ belong to an algebra, the series can be added, multiplied (using the Cauchy formula for product; note that for a fixed degree, computing the coefficients involves only a finite sum) or substituted one into another.

One can form equations among formal power series. The meaning is, of course, that the coefficients on each side should be the same. This is extremely useful in many areas of mathematics, notably combinatorics. See [Car95, Cos09] for more details on formal power series.
Many perturbation expansions in physics or in applied mathematics are based precisely into formulating the solutions of the equations of motion as formal power series and requiring that the equations of motion are satisfied in the sense of power series. Notably, the Lindstedt series were in standard use in astronomy even if they were only shown to converge for some frequencies in [Mos67].

2.3.2. Asymptotic expansions. For formal power series, a notion weaker that convergence of the series to a function is that the series is asymptotic to a function.

**Definition 8.** Given a set $D$ for which 0 is an adherent point, we say that a formal power series $\sum a_n \varepsilon^n$, with coefficients $a_n$ in a Banach space $X$, is an asymptotic expansion to a function $\phi : D \rightarrow X$ when for all $N \in \mathbb{N}_0$, there exists $C_N$ such that for any $\rho > 0$

$$\sup_{\varepsilon \in D, |\varepsilon| \leq \rho} \left\| \sum_{n=0}^{N} a_n \varepsilon^n - \phi(\varepsilon) \right\| \leq C_N \rho^{N+1}.$$ 

If the domain $D$ does not include any ball centered at zero, even if the function $\phi$ is analytic and bounded on $D$, this does not imply that the series converges.

Note that different functions may have the same asymptotic expansions. The Cauchy example

$$\phi(\varepsilon) = \exp(-\varepsilon^{-2})$$

has an identically zero asymptotic expansion on a domain

$$D_\delta = \{ \varepsilon : |\text{Arg}(\varepsilon)| < \delta \}$$

when $\delta < \pi/4$.

Note that the definition of asymptotic involves the domain $D$. A series may be asymptotic to a function in a domain but not in a larger domain. For example the zero series is asymptotic to the the Cauchy example 2.2 in the domains $D_\delta$ as in (2.3) when $\delta < \pi/4$, but not when $\delta > \pi/4$.

2.3.3. Gevrey formal expansions. Given a formal power series, even if it diverges, it is interesting to study how fast the coefficients grow. The following definition captures some speed of growth that is weaker than convergence, but which nevertheless appears naturally in many applied problems.

**Definition 9.** Let $\beta, \rho > 0$. We say that a power series expansion $f = \sum_{n=0}^{\infty} f_n(\theta)\varepsilon^n$, with $f_n \in A_\rho$, belongs to a Gevrey class $(\beta, \rho)$ if and only if there exist constants $C \geq 0, R \geq 0,$ and $n_0 \in \mathbb{N}$ such that

$$\|f_n\|_\rho \leq CR^\beta n^{\beta n}$$

for $n \geq n_0,$

and we denote $f \in G_\rho^\beta$.

Similarly, we say that a power series expansion $\mu = \sum_{n=0}^{\infty} \mu_n \varepsilon^n$, with $\mu_n \in \mathbb{C}_\rho$, belongs to a Gevrey class $\beta$ if and only if there exist constants $C \geq 0, R \geq 0,$ and $n_0 \in \mathbb{N}$ such that

$$|\mu_n| \leq CR^\beta n^{\beta n}$$

for $n \geq n_0,$

and we denote $\mu \in G^\beta$. 

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Remark 10. It is well known that (2.4) in definition 9 is equivalent to the inequality
\[ \| f_n \|_\rho \leq \hat{C} R^d (n!)^{\beta} \] for \( n \geq n_0 \)
which, in turn, implies the series \( \sum_{n=0}^{\infty} \frac{f(n)}{(n!)^{\beta}} \) converges in \( A_\rho \) with positive radius of convergence.

This remark makes a connection with the theory of Borel summability. If a series is Gevrey, under some extra conditions, the iterated Borel transform produces a function that is analytic in a sector and gives rise, by means of a variant of the Laplace transform, to a function to which the series is asymptotic. See [CGGG07, Cos09].

2.3.4. A property from number theory. In KAM theory, some number theoretical properties of frequencies play an important role.

Definition 11. Let \( \omega \in \mathbb{R}^d \), \( \tau > 0 \). We define the quantity \( \nu(\omega; \tau) \) as
\[ (\nu(\omega, \tau))^{-1} \equiv \sup_{k \in \mathbb{Z}^d \setminus \{0\}} |e^{2\pi ik \cdot \omega} - 1|^{-1} |k|^{-\tau}. \]

We say that \( \omega \) is Diophantine of class \( \tau \) and constant \( \nu = \nu(\omega, \tau) \), whenever \( \nu^{-1} < \infty \). We use the notation \( \omega \in D(\nu, \tau) \).

Note that if \( \omega \in D(\nu, \tau) \), then \( |e^{2\pi ik \cdot \omega} - 1| \geq \nu |k|^{-\tau} \).

2.4. Quasi-periodic orbits

A quasi-periodic sequence \( \{x_n\}_{n \in \mathbb{Z}} \) of frequency \( \omega \in \mathbb{R}^d \) in an Euclidean space is a sequence which can be expressed in terms of Fourier series.

\[ x_n = \sum_{k \in \mathbb{Z}^d} e^{2\pi i k \cdot \omega} \hat{x}_k = K(n\omega), \]

where \( K(\theta) = \sum_{k \in \mathbb{Z}^d} e^{2\pi i k \cdot \theta} \hat{x}_k \).

We can think of the function \( K \) as an embedding of the torus \( T^d \) into phase space. If \( \omega \) does not have any resonances (i.e. \( k \cdot \omega \neq 0 \) for \( k \in \mathbb{Z}^d \setminus \{0\} \), which can always be arranged by reducing \( d \) if there is one), then \( \{\omega n\}_{n \in \mathbb{Z}} \) is dense on the torus. The map \( K \) is often called the hull function.

If \( x_n \) is an orbit of a map, \( x_{n+1} = f(x_n) \) we see that \( K(n\omega + \omega) = f(K(n\omega)) \). Since \( \{\omega n\}_{n \in \mathbb{Z}} \) is dense, this is equivalent to
\[ K(\theta + \omega) = f(K(\theta)) \forall \theta \in T^d. \] (2.6)

Hence, we see that the set \( K(T^d) \), the image of the standard torus under the embedding \( K \) is invariant under \( f \). So, it is customary to describe quasi-periodic solutions as invariant tori.

The problem of given a map finding a quasi-periodic solution of frequency \( \omega \) can be formulated as finding an embedding \( K \) solving (2.6). The equation (2.6) will be our fundamental tool to characterize quasi-periodic orbits.
2.5. Set-up of the problem. The invariance equation

In this section, we describe informally the geometric set up and the geometric meaning of the formulation of our problem. The precise formulation of the main results of this paper theorems A and B will be presented in section 3.

We will be mainly concerned with an analytic family of maps \( f_{\varepsilon,\mu} : M \to M \), such that

\[ f_{\varepsilon,\mu}^* \Omega = \lambda(\varepsilon) \Omega, \]

where \( \varepsilon \in \mathbb{C} \) is a small parameter, \( \mu \in \Lambda \subseteq \mathbb{C}^d \) is an internal parameter (the drift parameter), and \( \lambda(\varepsilon) = 1 - \varepsilon^\alpha \).

A good example to keep in mind is the dissipative standard map presented in (1.1). Note that, for \( \varepsilon = 0 \) and for each \( \mu \), the maps \( f_{0,\mu} \) are symplectic because \( \lambda(0) = 1 \).

The main assumption in the main theorem, theorem B, is that the map \( f_{0,\mu} \) has an invariant torus in which the motion is a rotation of frequency \( \omega \) which is Diophantine (see definition 11).

Note that the drift parameter, \( \mu \), is chosen to guarantee the persistence of a quasi periodic orbit of a given frequency \( \omega \), so we also consider \( \mu = \mu_\epsilon \).

Following the discussion in section 2.4 and, in particular (2.6), we see that finding a quasi-periodic orbit for \( f_{\varepsilon,\mu} \) is equivalent to finding families of embeddings \( K_\varepsilon \) and families of parameters \( \mu_\epsilon \) in such a way that

\[ f_{\varepsilon,\mu_\epsilon} \circ K_\varepsilon(\theta) = K_\varepsilon(\theta + \omega). \tag{2.7} \]

Equation (2.7) should be interpreted as, given the family \( f_{\varepsilon,\mu} \) and the frequency \( \omega \) finding \( \mu_\epsilon, K_\varepsilon \). For this work, the sense in which (2.7) is meant to hold is the meaning of formal power series (the coefficients of \( \varepsilon^n \) on both sides of (2.7) are identical for all \( n \), as it is customary in the study of Lindstedt series).

Note that the equation (2.7) is highly underdetermined. If \( \mu_\epsilon, K_\varepsilon \) is a solution, changing \( \theta \) into \( \theta + \sigma \varepsilon \), we obtain that \( \mu_\epsilon, \tilde{K}_\varepsilon \) is also a solution where \( \tilde{K}_\varepsilon(\theta) = K_\varepsilon(\theta + \sigma \varepsilon) \). This change of variables has the physical meaning of choosing a change of origins in the torus.

2.6. Automatic reducibility

As it is noted in [CCdlL13], a very useful property of conformally symplectic systems is that solutions to equation (2.7) satisfy the so-called automatic reducibility, that is, in a neighborhood of an invariant torus, one can find a system of coordinates in which the linearization of the evolution has almost all the coefficients constant.

**Lemma 12.** Let \( f_\mu : M \to M \), such that, \( f_\mu^* \Omega = \lambda \Omega \), and \( K : T^d \to M \) such that

\[ f_\mu \circ K(\theta) = K(\theta + \omega) \]

with \( \omega \) an irrational vector. Denoting \( \mathcal{N} = (DK^\top DK)^{-1} \) and the \( 2d \times d \) matrix \( V(\theta) = J^{-1} \circ K(\theta)DK(\theta) \), then, the \( 2d \times 2d \) matrix

\[ M(\theta) = \begin{bmatrix} DK(\theta) & V(\theta) \end{bmatrix} \]

satisfies

\[ Df_\mu \circ K(\theta)M(\theta) = M(\theta + \omega) \begin{pmatrix} Id & S(\theta) \\ 0 & \lambda Id \end{pmatrix}, \tag{2.9} \]

where \( Id \in \mathbb{R}^{d \times d} \) and \( S(\theta) \) is an explicit algebraic expression involving \( DK, Df_\mu, J \circ K \), and, \( \mathcal{N} \); see (2.12).
The proof of lemma 12 is given in [CCdlL13]. The argument is as follows, taking derivative in equation (2.7) one has $Df_\mu \circ K_0(\theta)DK_0(\theta) = DK_0(\theta + \omega)$ which gives the first column in (2.9). The second column comes from the fact that the conformally symplectic property, $f_\mu^*\Omega = \lambda\Omega$, implies that the invariant torus given by equation (2.7) is Lagrangian. Then, using the conformally symplectic geometry the second column can be obtained.

**Remark 13.** As pointed out in [CCdlL13] if $K$ is an approximate solution of (2.7), that is,

$$f_\mu \circ K(\theta) - K(\theta + \omega) =: E(\theta)$$

(2.10)

then the relation (2.9) will hold with an error, $R$, that can be estimated in terms of the error, $E(\theta)$, of the invariance equation, that is

$$Df_\mu \circ K(\theta)M(\theta) = M(\theta + \omega)\begin{pmatrix} \Id & S(\theta) \\ 0 & \lambda\Id \end{pmatrix} + R(\theta),$$

(2.11)

with

$$S(\theta) = P(\theta + \omega)^T Df_\mu \circ K(\theta)J^{-1} \circ K(\theta)P(\theta)$$

$$P(\theta) = DK(\theta)N(\theta),$$

$$\Gamma(\theta) = DK(\theta)^T J^{-1} \circ K(\theta)DK(\theta).$$

(2.12)

Moreover, denoting $E_L(\theta) = DK(\theta)^T J \circ K(\theta)DK(\theta)$ as the error in the Lagrangian character of the torus, one has

$$R(\theta) = \left[DE(\theta) \begin{pmatrix} V(\theta + \omega)(\tilde{B}(\theta) - \lambda\Id) + DK(\theta + \omega)(\tilde{S}(\theta) - S(\theta)) \\ \Id \end{pmatrix}\right],$$

(2.13)

where

$$V(\theta) = J^{-1} \circ K(\theta)DK(\theta)N(\theta)$$

(2.14)

$$\tilde{B}(\theta) - \lambda\Id \equiv -E_L(\theta + \omega)\tilde{S}(\theta)$$

(2.15)

$$\tilde{S}(\theta) - S(\theta) \equiv -N(\theta + \omega)^T \Gamma(\theta + \omega)N(\theta + \omega)(\tilde{B}(\theta) - \lambda\Id).$$

(2.16)

We note that (2.15) and (2.16) can be considered as equations for $\tilde{B}$ and $\tilde{S}$, since they determine uniquely these quantities. As it is pointed out in [CCdlL13], this system have non vanishing diagonal terms and the upper diagonal terms are small. Thus, the system of equations (2.15) and (2.16) can be solved and one can bound the size of $\tilde{B} - \lambda\Id$ and $\tilde{S} - S$ by a constant times the size of $E_L$. Precise estimates are given in lemma 39. The derivation of the formulas in (2.12)–(2.14) can be found in [CCdlL13].

**Remark 14.** Observe that when considering $K_0$, $\mu_0$ satisfying (2.7) and a perturbation $K_\varepsilon$, $\mu_\varepsilon$ (which could be given in terms of formal power series), equation (2.11) is also satisfied by $K_\varepsilon$, $\mu_\varepsilon$ but with all the expressions depending on $\varepsilon$ (small enough), that is,
\[ Df_{\mu_\epsilon} \circ K_\epsilon(\theta)M_\epsilon(\theta + \omega) = M_\epsilon(\theta + \omega) \begin{pmatrix} 1 & \delta \theta \\ 0 & \lambda I \end{pmatrix} + R_\epsilon(\theta). \]

3. Statement of the main results

In this section we state the precise formulation of the main results, which give the Gevrey character of the formal expansions of the solutions to equation (2.7). First we introduce a normalization which guarantees the uniqueness of the solutions to equation (2.7).

3.1. Normalization and local uniqueness

The centerpiece of this work is the invariance equation

\[ f_{\epsilon, \mu_\epsilon} \circ K_\epsilon = K_\epsilon \circ T_\omega, \tag{3.1} \]

where \( T_\omega(\theta) = \theta + \omega \). Note that if \((K, \mu)\) is a solution of the invariant equation (3.1), then, for any \( \sigma \in \mathbb{T}^d \), \((K \circ T_\sigma, \mu)\) is also a solution of (3.1), due to the fact that \( K \circ T_\sigma \) parameterizes the same torus as \( K \). So, in order to get uniqueness it is necessary to impose a normalization condition.

**Definition 15.** Denote as \( K_0 \) a solution of (3.1) for \( \varepsilon = 0 \). We say that a torus with embedding \( K \) is normalized with respect to \( K_0 \) when

\[ \int_{\mathbb{T}^d} \left[ M_0^{-1}(\theta)(K(\theta) - K_0(\theta)) \right]_d d\theta = 0, \tag{3.2} \]

where the subscript \( d \) indicates that we take the first \( d \) rows of the \( 2d \times d \) matrix, and \( M_0 \) is constructed from \( K_0 \) as in (2.8).

We also recall the following result ([CCdlL13], proposition 26) which shows that this condition can be imposed without loss of generality for solutions that are close to one another.

**Proposition 16.** Let \( K_0, K \) be solutions of (3.1) and \( \|K - K_0\|_c \) be sufficiently small (with respect to quantities depending only on \( M \)-computed out of \( K_0 \)—and \( f \)). Then, there exists \( \sigma \in \mathbb{R}^d \), such that \( K^{(\sigma)} = K \circ T_\sigma \) satisfies (3.2). Furthermore,

\[ |\sigma| \leq C \|K - K_0\|_c, \]

where the constant \( C \) can be chosen to be as close to 1 as desired by assuming that \( f_{\epsilon, \mu}, K_0, \) and \( K_1 \) are twice differentiable, \( DK_0^\top DK_1 \) is invertible and \( \|K - K_0\|_c^0 \) is sufficiently small. The \( \sigma \) thus chosen is locally unique.

**Remark 17.** As it is noted in [CCdlL13] the normalization (3.2) works as well when \( K \) is only a formal solution. Then, assuming that \( K_0 \) is a solution of equation (3.1), the normalization condition (3.2) for a formal solution of (3.1) given as power series expansion \( \sum_{n=0}^{\infty} K_n(\theta) \varepsilon^n \) is equivalent to the conditions

\[ \int_{\mathbb{T}^d} [M_0^{-1}(\theta)K_n(\theta)]_d d\theta = 0 \tag{3.3} \]

for all \( n \geq 1 \).
3.2. Main theorems

We start this section stating our first main theorem, theorem A. The proof of theorem A is a consequence of theorem B. The proposition 57, given in the appendix A, shows that the hypothesis of theorem B are satisfied for maps of the form (3.4). Theorem B states the same results as theorem A but in a more general setting. In section 3.3 we state theorem C, which establishes that the formal expansions in theorem A are also asymptotic to geometric objects.

**Theorem A.** Let \( \omega \in \mathcal{D}(\nu, \tau) \). Consider the map \( f : \mathbb{T} \times \mathbb{R} \rightarrow \mathbb{T} \times \mathbb{R} \) given by

\[
f_{\varepsilon, \mu}(x, y) = (x + \lambda(\varepsilon)y + \mu \varepsilon - \varepsilon V'(x), \lambda(\varepsilon)y + \mu \varepsilon - \varepsilon V'(x)),
\]

where \( \lambda(\varepsilon) = 1 - \varepsilon^\alpha, \alpha \in \mathbb{N}, V(x) \) is a trigonometric polynomial, \( \mu, \varepsilon \in \mathbb{C} \), and \( \varepsilon \in \mathbb{C} \). Then, there exists \( \rho_0 > 0 \) such that the following holds

(a) There exist formal power series expansions \( K_n^{(x)} = \sum_{j=0}^\infty K_{\omega}^{(x)} \) and \( \mu_\varepsilon^{(x)} = \sum_{j=0}^\infty \mu_\varepsilon^{(x)} \) satisfying \( f_{\varepsilon, \mu} \circ K = K(\theta + \omega) \) in the sense of formal power series. More precisely, defining \( K_{\omega}^{(x)} = \sum_{j=0}^N K_{\omega}^{(x)} \) and \( \mu_\varepsilon^{(x)} = \sum_{j=0}^N \mu_\varepsilon^{(x)} \) for any \( N \in \mathbb{N} \) we have

\[
\left\| f_{\varepsilon, \mu}^{(x)} \circ K_{\omega}^{(x)} - K_{\omega}^{(x)} \circ T_\omega \right\| \leq C_N |\varepsilon|^{N+1},
\]

where \( C_N > 0 \). Moreover, if the \( K_j \)'s satisfy the normalization condition (3.3), then the expansions \( K_n^{(x)} \), \( \mu_\varepsilon^{(x)} \) are unique.

(b) The unique formal power series expansions, \( K_n^{(x)} \) and \( \mu_\varepsilon^{(x)} \), satisfying (3.5) and the normalization (3.3) are such that \( K_n^{(x)} \in G_{\rho_0/2}^{(x)} \) and \( \mu_\varepsilon^{(x)} \in G^{(x)}_{\tau/\alpha} \), i.e. there exists constants \( L, F, N_0 \) such that

\[
\|K_n\|_{\mu} \leq LF^n n^{(\tau/\alpha)n} \quad \text{and} \quad |\mu_n| \leq LF^n n^{(\tau/\alpha)n} \quad \text{for any } n > N_0.
\]

(3.6)

**Remark 18.** It is instructive to compare the results in theorem A with the numerical explorations of [BC19] (see also [BC21]). In the case that \( \lambda(\varepsilon) = 1 - \varepsilon^3 \) and \( \omega \) is the golden mean, theorem A gives that the expansion satisfies the Gevrey bounds with exponent 1/3. Of course, theorem A gives only an upper bound and lower exponents could also be true. The numerical results in [BC19] and [BC21] lead to the conjecture that the expansion \( \sum K_n \varepsilon^n \) has some well defined asymptotics

\[
\|K_n\|_{\mu}^{1/n} \lesssim Cn^\tau.
\]

The numerical values of the Gevrey exponent for several values of the frequency, Diophantine with \( \tau = 1 \), were computed in [BC21] and the largest one found was 0.3.

The asymptotics (3.7) is compatible with the results in theorem A and indeed, the numerical values of the exponents are close to the rigorous upper bounds in some cases.

We call attention that [BC19] contained an unfortunate typo and the results attributed there to \( \|K_n\|^{1/n} \) are actually results for \( \|nK_n\|^{1/n} \), this is corrected in [BC21]. The paper [BC21] also presents several other patterns in the series (refined versions of (3.7) including oscillations of period 3, studies for other Diophantine numbers, etc). We hope that the method presented in this paper can lead to studies of these phenomena, hitherto discovered only through numerical exploration.
We think that the argument in theorem A can be adapted to remove the hypothesis of \( V(0) \) being a trigonometric polynomial (it may suffice that the iterative step involves an extra loss of domain in the angle \( \theta \)). Since the method of proof is rather novel, we decided to follow the advice “Premature optimization is the root of all evil” [Knu98], and present the argument in its simplest form so that it could, perhaps, be applied to other problems.

For the sake of completeness, before stating the main theorem B we will state a theorem in [CCdlL17] which assures the existence of formal power series expansions satisfying (3.1) up to any order for conformally symplectic systems. We note that the family of maps (3.4) considered in theorem A satisfies the hypothesis of theorem 19.

**Theorem 19 ([CCdlL17], theorem 12).** Let \( \mathcal{M} \equiv \mathbb{T}^d \times B \) with \( B \subseteq \mathbb{R}^d \) an open, simply connected domain with smooth boundary; \( \mathcal{M} \) is endowed with an analytic symplectic form \( \Omega \).

Let \( \omega \in \mathcal{D}(\tau, \nu) \) and consider a family \( f_{\varepsilon, \mu} \) of conformally symplectic mappings that satisfy

\[
f_{\varepsilon, \mu}^* \Omega = \lambda(\varepsilon) \Omega,
\]

with \( \mu \in \Lambda \), with \( \Lambda \subseteq \mathbb{C}^d \) an open set, \( \lambda(\varepsilon) = 1 - \varepsilon^\alpha \), \( \alpha \in \mathbb{N} \) and \( \varepsilon \in \mathbb{C} \).

Assume that \( f_{\varepsilon, \mu} \) depends analytically on \( (\varepsilon, \mu) \) for \( \varepsilon \) near to 0 and \( \mu \in \Lambda \). Assume also that for \( \varepsilon = 0 \) the family of maps \( f_{0, \mu} \) is symplectic and that for some value \( \mu_0 \in \Lambda \) the map \( f_{0, \mu_0} \) admits a Lagrangian invariant torus, namely we can find an analytic embedding \( K_0 \in \mathcal{A}_\mu(\mathbb{R}^d, \mathcal{M}) \), for some \( \rho > 0 \), such that

\[
f_{0, \mu_0} \circ K_0 = K_0 \circ T_\omega.
\]

Furthermore, assume that the torus \( K_0 \) satisfies the following hypothesis:

\[\text{HND.} \text{ Let the following non-degeneracy condition be satisfied:}\]

\[
det \begin{pmatrix}
S_0 & \tilde{S}_0(B_{0b})^0 + \tilde{A}_{01} \\
0 & \tilde{A}_{02}
\end{pmatrix} \neq 0
\]

where the \( d \times d \) matrix \( S_0 \) is defined as

\[
S_0(\theta) \equiv N_0(\theta + \omega)^\top D\theta(\theta + \omega)^\top Df_{0,0} \circ K_0(\theta)J^{-1} \circ K_0(\theta)DK_0(\theta)N_0(\theta)
\]

\[\quad - N_0(\theta + \omega)^\top DK_0(\theta + \omega)^\top J^{-1} \circ K_0(\theta + \omega)DK_0(\theta + \omega)N_0(\theta + \omega)\]

with \( N = (DK_0 \circ DK_0)^{-1} \), the \( d \times d \) matrices \( \tilde{A}_{01}, \tilde{A}_{02} \) denote the first \( d \) and the last \( d \) rows of the \( 2d \times d \) matrix \( \tilde{A}_0 = (M_0 \circ T_\omega)^{-1} (D_{\rho}f_{0,\mu_0} \circ K_0) \), where \( M_0 \) is as in (2.8), \( (B_{0b})^0 \) is the solution (with zero average) of the cohomology equation \( (B_{0b})^0 = - (\tilde{A}_{02})^0 \), where \( (B_{0b})^0 \equiv B_{0b} - \overline{B_{0b}} \) and the overline denotes the average.

Then, we have the following

(a) There exist a formal power series expansions \( K_{\varepsilon}^{[\infty]} = \sum_{j=0}^{\infty} K_j \varepsilon^j \) and \( \mu_{\varepsilon}^{[\infty]} = \sum_{j=0}^{\infty} \mu_j \varepsilon^j \) satisfying (3.9) in the sense of formal power series. More precisely, defining \( K_{\varepsilon}^{[N]} = \sum_{j=0}^{N} K_j \varepsilon^j \) and \( \mu_{\varepsilon}^{[N]} = \sum_{j=0}^{N} \mu_j \varepsilon^j \) for any \( N \in \mathbb{N} \) and \( \rho > 0 \), we have

\[
\left\| f_{\varepsilon, \mu}^{[N]} \circ K_{\varepsilon}^{[N]} - K_{\varepsilon}^{[N]} \circ T_\omega \right\|_{\rho_0} \leq C_N |\varepsilon|^{N+1}.
\]

(3.10)
for some $0 < \rho_0 < \rho$ and $C_N > 0$.

Moreover, if we require the $K_j$’s satisfy the normalization condition (3.3), then the expansions $K_1^{[\infty]}$, $\mu_1^{[\infty]}$ are unique.

Note that theorem 19 does not assume that the case $\varepsilon = 0$ is an integrable system, as it is the case for the map (3.4), it suffices that the case $\varepsilon = 0$ is a symplectic system with a KAM torus.

**Remark 20.** Denoting

$$E^N_\varepsilon(\theta) = f_{\varepsilon, \mu_j, \omega} \circ K^1_{\leq N}(\theta) - K^1_{\leq N}(\theta + \omega)$$

then (3.10) can be written as

$$\left\| E^N_\varepsilon \right\|_{\mu_0} \leq C_N |\varepsilon|^{N+1}.$$  \hspace{1cm} (3.11)

According to the notation introduced earlier, this means that $E^N_\varepsilon \sim O(|\varepsilon|^{N+1})$ or $E^N_\varepsilon = \sum_{j=N+1}^{\infty} E_j \varepsilon^j$ for $\varepsilon$ small enough. We denote

$$E^{N,2N}_\varepsilon = \sum_{j=N+1}^{2N} E_j \varepsilon^j$$

the truncated series.

The following theorem, theorem B, can be considered as an improvement of theorem 19 in the sense that it gives Gevrey bounds for the coefficients $K_j$, $\mu_j$ of the unique (under normalization) formal power series expansions $K_1^{[\infty]}$, $\mu_1^{[\infty]}$.

**Theorem B.** Assume the hypotheses of theorem 19. Assume also that for any $\varepsilon$, small enough, and for any $N \in \mathbb{N}$ we have:

**HTP1.** $\tilde{E}_{1,2}^{N,2N}$, $\tilde{A}_{1,2}^{N}$ are trigonometric polynomials in $\theta$ of degree at most $aN$, $a \in \mathbb{N}$. Where $\tilde{E}_{1,2}^{N,2N}$, $\tilde{A}_{1,2}^{N}$ denote the $d \times 1$ and $d \times d$ matrices, respectively, given by taking the last $d$ rows of the $2d \times 1$ matrix $\tilde{E}_{1,2}^{N,2N} = (M^1_{\varepsilon, \mu, \omega} \circ T_\omega)^{-1} E_{1,2}^{N,2N}$ and the $2d \times d$ matrix $\tilde{A}_{1,2}^{N} = (M^1_{\varepsilon, \mu, \omega} \circ T_\omega)^{-1} D_{\varepsilon, \mu, \omega} \circ K_2^{[\leq N]}$, respectively. $M^1_{\varepsilon, \mu, \omega}$ is as in (2.8) constructed from $K_1^{[\leq N]}$.

**HTP2.** The $d \times d$ matrix

$$\tilde{E}_{1,2}^{N}(\theta) \equiv DK_2^{[\leq N]}(\theta + \omega) \circ K_2^{[\leq N]}(\theta + \omega) D_{\varepsilon, \mu, \omega} \circ K_2^{[\leq N]}(\theta + \omega)$$

$$- D(f_{\varepsilon, \mu, \omega} \circ K_2^{[\leq N]}(\theta)) \circ J \circ (f_{\varepsilon, \mu, \omega} \circ K_2^{[\leq N]}(\theta))$$

$$\times D(f_{\varepsilon, \mu, \omega} \circ K_2^{[\leq N]}(\theta))$$  \hspace{1cm} (3.13)

is a trigonometric polynomial of degree at most $aN$.

Then, the unique formal power series expansions, $K_1^{[\infty]}$ and $\mu_1^{[\infty]}$, satisfying (3.10) and (3.3) are such that $K_1^{[\infty]} \in \mathcal{G}^{(r)/(\alpha)}_{0/2}$ and $\mu_1^{[\infty]} \in \mathcal{G}^{(r)/(\alpha)}$, i.e. there exists constants $L$, $F$, $N_0$ such that

$$\|K_n\|_{2N} \leq LF^n n^{(r)/(\alpha)n}$$

and

$$|\mu_n| \leq LF^n n^{(r)/(\alpha)n}$$

for any $n > N_0$.  \hspace{1cm} (3.14)
The proof of theorem B, given in section 6.2, is done by means of a Newton like method which acts on finite powers series expansions \((K_\varepsilon \subseteq \mathbb{N}, \mu_\varepsilon \subseteq \mathbb{N})\), this method is described in the next section. We emphasize that this quasi Newton method takes advantage of the conformally symplectic property (see definition 1) that maps like (3.4) satisfy.

We also point out that hypothesis HTP1 and HTP2 are natural for the maps considered in theorem A. The verification of these hypothesis for the dissipative standard map is described in detail in proposition 57 of the appendix A. In the general setting in which theorem B is stated, the hypothesis HTP1 and HTP2 are needed to be able to get estimates, in balls with center at the origin, for the solutions of the linear equations of the quasi Newton method.

### 3.3. Asymptotic estimates for invariance functions

The formal power series studied in this paper are asymptotic expansions of functions \(K_\varepsilon, \mu_\varepsilon\) constructed in [CCdlL17]. The functions \(K_\varepsilon, \mu_\varepsilon\) are determined by the condition that they satisfy the invariance equation (3.1) and the normalization (3.3). In this section we argue that the same method we use to prove the Gevrey estimates also shows that the formal series defined here are asymptotic to the functions \(K_\varepsilon, \mu_\varepsilon\) with very strong estimates in the remainder, see theorem C.

We emphasize that the functions \(K_\varepsilon, \mu_\varepsilon\) are not constructed out of the asymptotic expansions by complex analysis methods (Borel summation, resummation of series). They are obtained from the requirement that they satisfy the invariance equation (3.1) and the normalization (3.3). It is an interesting open question whether some resummation of the asymptotic expansions studied here can produce the functions \(K_\varepsilon, \mu_\varepsilon\).

The domain of definition of the functions \(K_\varepsilon, \mu_\varepsilon\) is rather subtle. In [CCdlL17], it is proved that the domain of definition of \(K_\varepsilon, \mu_\varepsilon\) contains a set \(G\) obtained by removing sequence of balls that are dense on curves converging to the origin, in fact, it is rigorously showed that \(G\) is a lower bound on the analyticity domain of the functions \(K_\varepsilon, \mu_\varepsilon\). We also point out that the set \(G\) does not contain any ball centered at the origin. Indeed, the set \(G\) does not contain any sector centered at the origin of width bigger than \(\pi/\alpha\) (but it does contain a bounded sector centered on \(\mathbb{R}^+\) with opening \(<\pi/\alpha\), thus the width of the domain is not enough to apply many methods of complex analysis related to Phragmén–Lindelöf theory. In the other direction, the paper [CCdlL17] contains arguments showing that for generic perturbations one should not expect that the domain of analyticity contains the excluded balls (if the perturbation happens to be identically zero one indeed obtains a larger domain). The paper [BC19] studies numerically the maximal domain of definition of the functions \(K_\varepsilon, \mu_\varepsilon\) for the map (3.4) using a variety of methods including Padé summation and continuation methods. Indeed [BC19] conjectured that the series were Gevrey and this was an important motivation for this paper.

The set \(G\) is determined by asking that \(\lambda(\varepsilon)\) satisfies a Diophantine condition with respect to \(\omega\), more precisely, given \(\lambda \in \mathbb{C}\) we define

\[
\tilde{\nu} = \tilde{\nu}(\lambda; \omega, \tau) \equiv \sup_{k \in \mathbb{Z} \setminus \{0\}} |e^{2\pi ik\omega} - \lambda|^{-1}|k|^{-\tau}
\]  

and recalling that \(\lambda(\varepsilon) = 1 - \varepsilon^\alpha\)

\[
G = G(A; \omega, \tau, N) = \left\{ \varepsilon \in \mathbb{C} : \tilde{\nu}(\lambda(\varepsilon); \omega, \tau) |\lambda(\varepsilon) - 1|^{N+1} \leq A \right\} = \left\{ \varepsilon \in \mathbb{C} : \tilde{\nu}(\lambda(\varepsilon); \omega, \tau) \leq A |\varepsilon|^{-\alpha(N+1)} \right\}
\]

**Remark 21.** The meaning of the set \(G\) is that it is the set of \(\varepsilon \in \mathbb{C}\) for which the Diophantine constant \(\tilde{\nu}(\lambda(\varepsilon); \omega, \tau)\) is bounded by a power of \(|\varepsilon|\).
Remark 22. We note that, for fixed $\omega$, the function $\tilde{\nu} = \tilde{\nu}(\lambda(\varepsilon); \omega, \tau)$ is a lower semi-continuous function of $\varepsilon$, since it is the supremum of continuous functions. We also note that if $\omega \in \mathcal{D}(\nu, \tau)$, then $\tilde{\nu}(\lambda(0); \omega, \tau) = \nu^{-1} > 0$.

Remark 23. We point out that the set $\mathcal{G}$ does not contain any ball $B_k$ with center in $\hat{\varepsilon}_k$ such that $\Lambda(\hat{\varepsilon}_k) = e^{2\pi k \omega}$. Since $\lambda(\varepsilon) = 1 - \varepsilon^\alpha$, it is easy to construct a sequence of complex numbers $\varepsilon_k$ such that $\Lambda(\varepsilon_k) = e^{2\pi k \omega} = \varepsilon_k \to 0$, therefore $\mathcal{G}$ does not contain any ball centered at the origin. Indeed, in [CCdlL17], proposition 24, it is shown that ‘generic’ perturbations cannot lead to a formal expansions of the functions $K_\varepsilon$, $\mu_\varepsilon$ around $\varepsilon_k$. Thus the functions cannot be analytic at any $\varepsilon_k$.

The argument in [CCdlL17] does not consider the case of non-generic perturbations (like trigonometric polynomials), but it seems that the argument can be adapted.

The basic idea to prove the existence of the functions $K_\varepsilon$, $\mu_\varepsilon$ is as follows: the formal power expansions produces a sequence of polynomials which satisfy the invariance equation (3.1) rather approximately in a ball. In the intersection of the ball with the set $\mathcal{G}$, we can apply the a-posteriori theorem, theorem 14 in [CCdlL17], and obtain a true solution of (3.1). Of course, the detailed implementation requires taking into account several other issues such as the absence of monodromy.

In this paper we will use a very similar technique. As a byproduct of the estimates used in the proof of theorem B, we obtain that some truncations of the formal expansion satisfy the invariance equation up to a very small error in appropriate balls with sufficiently small radius. Then, in the intersection of the balls with the set $\mathcal{G}$ we will be able to apply theorem 14 in [CCdlL17].

More precisely we have:

**Theorem C.** Define $\hat{\gamma}_h = (2^{-1} \nu)^{1/\alpha}(a2^h N_0)^{-\tau/\alpha}$, with $h \in \mathbb{N}_0$. Assume the hypotheses of theorem B and that for any $0 < \delta < \rho_0$ one has that for all $\varepsilon \in \mathcal{G}$ such that $|\varepsilon| \leq \hat{\gamma}_0$,

$$
\|E_{\varepsilon,\omega}^{(\nu)}\|_{\mathcal{M}} \leq C(\nu^{-1} \tilde{\nu}(\lambda(\varepsilon), \omega, \tau))^2 \delta^{\nu + 1}.
$$

(3.17)

Let $\sum K_{\varepsilon} \in \ell^1$ and $\sum \mu_{\varepsilon} \in \ell^1$ the Gevrey formal expansion given by theorem B, and $n \in (2^h N_0, 2^{h+1} N_0] \cap \mathbb{N}$. Then, for any $\varepsilon \in \mathcal{G}$ such that $|\varepsilon| \leq \hat{\gamma}_{h+1}$ there exist $\mu_\varepsilon, K_\varepsilon$ such that $f_{\varepsilon, \lambda}, K_\varepsilon = K_\varepsilon \circ \tau_\omega$. Moreover, $K_\varepsilon$, $\mu_\varepsilon$ satisfy the inequalities

$$
\|K_\varepsilon - \sum_{j=0}^{n} K_{\varepsilon} \|_{\mathcal{M}} \leq \tilde{c}(CD)^{b} B^{2} r^{2} r^{2} N_0
$$

(3.18)

$$
|\mu_\varepsilon - \sum_{j=0}^{n} \mu_{\varepsilon} | \leq \tilde{c}(CD)^{b} B^{2} r^{2} r^{2} N_0,
$$

(3.19)
where \( \tilde{c} \) and \( C \) are uniform constants and 
\[
D = \nu^{-2}(\alpha N_0)^{2r} \rho_0^{-(2r+6d)} 2^{-(4r+12d)}, \quad r = 2^{-\tau/\alpha} \text{ and } B = 2^{\nu+6d}.
\]

Note that (3.18) can be understood as having super-exponentially small errors in domains decreasing exponentially fast. It is also important to note that almost all constants in (3.18) are given explicitly. We note that the hypotheses of theorem C are satisfied by the family of maps (3.4), thus the same result is obtained for the family of maps considered in theorem A. The proof of theorem C is given in section 6.3.

**Remark 24.** It is interesting to compare the inequalities (3.18) and (3.19) with the stronger notion of Gevrey asymptotic expansion in sectors. Given a function \( g \) (possibly taking values in a Banach space) analytic in an angular sector \( S \) and a formal power series \( \tilde{g} = \sum g_i \varepsilon^i \), one says that \( \tilde{g} \) is a \( \sigma \)-Gevrey asymptotic expansion of \( g \) if for every closed subsector \( \bar{S} \subset S \) there exist \( c_0, c_1 > 0 \) such that for any \( n > 0 \)
\[
\| g_{\varepsilon} - \sum_{i=0}^{n-1} g_i \varepsilon^i \| \leq c_0 c_1^n \sigma^n \varepsilon^n. \tag{3.20}
\]
If one considers \( |\varepsilon| \lesssim n^{-\tau/\alpha} \), as in theorem C, then (3.20), with \( \sigma = \tau/\alpha \), implies that
\[
\| g_{\varepsilon} - \sum_{i=0}^{n-1} g_i \varepsilon^i \| \leq c_0 c_1^n. \tag{3.21}
\]
Note that the estimates given in theorem C are estimates of the form (3.21), but in domains in \( \varepsilon \) that decreases as a power of \( n \).

It would be interesting to know if some quantitative Harnack inequalities could allow to obtain control in uniform domains from the control in the small domains obtained.

**Remark 25.** The use of resummation methods usually allows one to obtain estimates of the form (3.20). As it is mentioned above, the functions \( K_\varepsilon \) and \( \mu_\varepsilon \) are not constructed out of the asymptotic expansion by complex analysis methods, they are obtained from the requirement that they satisfy the invariance equation (3.1) and the normalization (3.3). We think that, once that the series is known to be Gevrey, resummation methods might be able to be applied to obtain functions \( \tilde{K}_\varepsilon \) and \( \tilde{\mu}_\varepsilon \) satisfying (3.20). However, even if it is possible to construct such a function \( \tilde{K}_\varepsilon \) by using resummation techniques, we are not sure if this function would be the parameterization of a quasi periodic orbit, that is, we do not know if \( \tilde{K}_\varepsilon \) would satisfy the invariance equation \( f_{\varepsilon,\mu_\varepsilon} \circ K_\varepsilon = K_\varepsilon \circ T_\omega \). The construction of a quasi periodic orbit through the use of resummation techniques would be a valuable result.

4. Iterative step of the quasi Newton method

The KAM procedure for the proof of theorem B is based on the application of a quasi Newton method, which is described in section 4.2. Before describing this procedure we introduce two types of cohomology equations that allow us to solve the linear equations, and obtain estimates, of the modified Newton method. The estimates for each step of the method will be given in section 5.
4.1. Estimates for some cohomology equations

The iterative step described in section 4.2 depends on the solution of two cohomology equations. The first equation, (4.1), is very standard in KAM theory. The estimate given in lemma 26 is well known for the experts in KAM theory, we have decided to include a proof here for the sake of completeness. The second type of cohomology equation we consider, (4.3), it is more complicated to study due to the appearance of the factor \( \lambda(\epsilon) = 1 - \epsilon^\alpha \). This factor introduces some restrictions in the set of parameters, \( \epsilon \), for which we are able to obtain estimates.

4.1.1. Standard cohomology equation. The first cohomology equation we deal with is the following

\[
\phi_\epsilon(\theta) - \phi_\epsilon(\theta + \omega) = \eta_\epsilon(\theta). \tag{4.1}
\]

Lemma 26 below, gives sufficient conditions to solve equation (4.1) and to obtain estimates of its solutions. This estimates are very standard in KAM theory.

**Lemma 26.** Let \( \omega \in D(\nu, \tau) \). Assume that \( \eta \in A_{\rho, r} \) is such that \( \int_\mathbb{T} \eta_\epsilon(\theta) d\theta = 0 \). Then, we can find a unique solution of (4.1), \( \phi_\epsilon \), that satisfies \( \int_\mathbb{T} \phi_\epsilon(\theta) d\theta = 0 \). Moreover, if \( 0 < \delta \leq \rho \), then \( \phi \in A_{\rho - \delta, r} \) and

\[
\| \phi \|_{\rho - \delta, r} \leq C \nu^{-1}\delta^{-(\tau + d)} \| \eta \|_{\rho, r}.
\]

with \( C = C(d) \). Furthermore, \( \eta_\epsilon \sim O(\| \epsilon \|^k) \) implies \( \phi_\epsilon \sim O(\| \epsilon \|^k) \).

**Proof.** Expanding (4.1) in Fourier series, the zero-mean-value solution is given by \( \phi_\epsilon(\theta) = \sum_{k \in \mathbb{Z} \setminus \{0\}} \hat{\eta}_k(\epsilon) \frac{e^{2\pi ik \cdot \theta}}{1 - e^{2\pi i k \omega}} \). Then, using Cauchy estimates one obtains

\[
\| \phi_\epsilon \|_{\rho - \delta, r} \leq \sum_{k \in \mathbb{Z} \setminus \{0\}} \left| \hat{\eta}_k(\epsilon) \right| \left\| \frac{e^{2\pi ik \cdot \theta}}{1 - e^{2\pi i k \omega}} \right\|_{\rho - \delta, r}
\]

\[
\leq \sum_{k \in \mathbb{Z} \setminus \{0\}} \nu^{-1}|k| \| \eta_\epsilon \|_\rho \| e^{2\pi ik \cdot \theta} \|_{\rho - \delta, r}
\]

\[
\leq C \nu^{-1} \| \eta_\epsilon \|_\rho \sum_{j \in \mathbb{N}} (\rho - \delta)^{-d - 1} e^{-2\pi \delta j}
\]

\[
\leq C \nu^{-1} e^{-(\tau + d)} \| \eta_\epsilon \|_\rho. \tag{4.2}
\]

The last line gives \( \phi_\epsilon \sim O(\| \epsilon \|^k) \) if \( \eta_\epsilon \sim O(\| \epsilon \|^k) \) and taking supremum over \( \epsilon \) the result is proved. \( \square \)

**Remark 27.** Equation (4.1) appears very often in KAM theory. When \( \epsilon \in \mathbb{R} \), the paper [Rüs75] contains estimates with a better exponent on \( \delta \). That is, in the same situation of lemma 26, when \( \epsilon \in \mathbb{R} \), one can get \( \| \phi_\epsilon \|_{\rho - \delta} \leq C \nu^{-\tau} \| \eta_\epsilon \|_\rho \).
4.1.2. Parametric cohomology equation. The second cohomology equation we are interested in is an equation for \( \varphi_\varepsilon : \mathbb{T}^d \to \mathbb{C} \), of the form
\[
\lambda(\varepsilon) \varphi_\varepsilon(\theta) - \varphi_\varepsilon(\theta + \omega) = \eta_\varepsilon(\theta), \tag{4.3}
\]
where \( \eta_\varepsilon : \mathbb{T}^d \to \mathbb{C} \) and \( \omega \in \mathbb{R}^d \) are given, \( \varepsilon \) fixed.

Note that, as it is seen in lemma 29, solve equation (4.3) presents a small divisors problem. In this case the small divisors depend on the variable \( \varepsilon \), that is, equation (4.3) is not expected to have a solution when \( \lambda(\varepsilon) = e^{2\pi i k \omega} \). One approach that has been used to deal with the small divisors in equation (4.3) (see [CCdlL17]) requires to remove a set from the complex plane, \( \varepsilon \in \mathbb{C} \), where the denominators \( \lambda(\varepsilon) = e^{2\pi i k \omega} \) are small. This gives rise to a set with a complicated structure, \( \mathcal{G} \subset \mathbb{C} \), of parameters, \( \varepsilon \), in which it is possible to find a solution, and estimates, of equation (4.3). One of the properties of the set \( \mathcal{G} \) described in [CCdlL17], is that it does not contain any ball with center at the origin. This property is one of the reasons for which we follow a different approach to deal with equation (4.3): to prove the Gevrey estimates in theorem B we rely heavily on being able to obtain estimates of (4.3) for \( \varepsilon \) in a ball centered at the origin.

The following two lemmas allow us to obtain estimates in balls centered at \( \varepsilon = 0 \) for the solution, \( \varphi_\varepsilon \), of equation (4.3) whenever \( \eta_\varepsilon \) is a trigonometric polynomial. If the degree of the trig polynomial, \( \eta_\varepsilon \), is \( aN \), lemma 28 gives a relation between this degree and a domain in which the solution, \( \varphi_\varepsilon \), of (4.3) will be analytic in \( \varepsilon \).

Note that the requirement of hypotheses HTP1 and HTP2 in theorem B is due to the fact that the quantities given in these hypothesis will be the right-hand side of equations of the form (4.3).

**Lemma 28.** Let \( \omega \in \mathcal{D}(\nu, \tau) \), \( \lambda(\varepsilon) = 1 - \varepsilon^\alpha \), \( \alpha \geq 1 \), and \( a, N \in \mathbb{N} \). If \( |\varepsilon| \leq \left( \frac{\nu}{2} \right)^{1/\alpha} \frac{1}{(aN)^r} \), then, for \( |k| \leq aN \) we have
\[
|\lambda(\varepsilon) - e^{2\pi i k \omega}| \geq \frac{\nu}{2} \frac{1}{(aN)^r}.
\]

**Proof.**
\[
|e^{2\pi i k \omega} - \lambda(\varepsilon)| \geq |e^{2\pi i k \omega} - 1| - |1 - \lambda(\varepsilon)| \geq \frac{\nu}{|k|^r} - |\varepsilon|^\alpha
\]
\[
\geq \frac{\nu}{(aN)^r} - \frac{\nu}{2(aN)^r} = \frac{\nu}{2} \frac{1}{(aN)^r}.
\]

**Lemma 29.** Let \( \lambda(\varepsilon) = 1 - \varepsilon^\alpha \), \( \alpha \geq 1 \), \( \omega \in \mathcal{D}(\nu, \tau) \); \( a, N \in \mathbb{N} \), and define
\[
\gamma_N = \left( \frac{\nu}{2} \right)^{1/\alpha} \frac{1}{(aN)^r}.
\]

Let \( \eta \in \mathcal{A}_{\rho, \gamma_N} \) such that \( \int_{\mathbb{T}} \eta(\theta) d\theta = 0 \) and assume that, for any \( \varepsilon \), \( \eta_\varepsilon(\theta) \) is a trigonometric polynomial of degree \( aN \) in \( \theta \). Then, for any \( |\varepsilon| \leq \gamma_N \) equation (4.3) has a unique solution, \( \varphi_\varepsilon(\theta) \), such that \( \int_{\mathbb{T}} \varphi_\varepsilon(\theta) d\theta = 0 \). Furthermore, if \( 0 < \delta \leq \rho \), then \( \varphi \in \mathcal{A}_{\rho - \delta, \gamma_N} \) and
\[
\|\varphi\|_{\rho - \delta, \gamma_N} \leq C \nu^{-1}(aN)^r \delta^{-d} \|\eta\|_{\rho, \gamma_N}.
\]
Moreover, if \( \eta_\varepsilon \sim \mathcal{O}(|\varepsilon|^d) \), then \( \varphi_\varepsilon \sim \mathcal{O}(|\varepsilon|^d) \).
Proof. Expanding $\eta_\varepsilon$ in Fourier series as $\eta_\varepsilon(\theta) = \sum_{0 < |k| = aN} \hat{\eta}_k(\varepsilon)e^{2\pi ik\theta}$, the zero-mean-value solution to $(4.3)$ is given by

$$\varphi_\varepsilon(\theta) = \sum_{0 < |k| = aN} \frac{\hat{\eta}_k(\varepsilon)}{N(\varepsilon)} e^{2\pi ik\theta}.$$  

Using lemma 28 and Cauchy estimates, one obtains that for any $|\varepsilon| \leq \gamma_N$

$$\|\varphi_\varepsilon\|_{\mu-\delta} \leq \sum_{0 < |k| = aN} \frac{|\hat{\eta}_k(\varepsilon)|}{N(\varepsilon)} \|e^{2\pi ik\theta}\|_{\mu-\delta} \leq 2(aN)^{-\nu^{-1}} \sum_{0 < |k| \leq aN} |\hat{\eta}_k(\varepsilon)|^{2|\rho(k)|^{(\mu-\delta)}}$$

$$\leq 2(aN)^{-\nu^{-1}} \sum_{0 < |k| \leq aN} \|\eta_\varepsilon\|_\rho e^{2|\rho(k)|^{(\mu-\delta)}} \leq 2(aN)^{-\nu^{-1}} \|\eta_\varepsilon\|_\rho \sum_{j=1}^{aN} j^{\delta-1} e^{-2\pi j\varepsilon} \leq C\nu^{-1}(aN)^{-\delta^{-d}} \|\eta_\varepsilon\|_{\mu^{-N}}.$$  

Thus, $\|\varphi_\varepsilon\|_{\mu^{-N},\eta_\varepsilon} \leq C\nu^{-1}(aN)^{-\delta^{-d}} \|\eta_\varepsilon\|_{\mu^{-N}}$. The last claim comes from $(4.4)$, that is $\varphi_\varepsilon \sim O(|\varepsilon|^k)$ if $\eta_\varepsilon \sim O(|\varepsilon|^k)$. \hfill $\square$

### 4.2. Formulation of the quasi Newton method

Every step of the quasi Newton method starts with a solution of equation $(3.1)$ up to order $\varepsilon^N$. That is, assume that

$$K_\varepsilon^{(N)}(\theta) = \sum_{n=0}^{N} K_n(\theta)e^n, \quad \mu_\varepsilon^{(N)} = \sum_{n=0}^{N} \mu_ne^n$$

satisfy the normalization $(3.3)$ and

$$f_{c,\mu;\varepsilon} \circ K_\varepsilon^{(N)}(\theta) = K_\varepsilon^{(N)}(\theta + \omega) = E_\varepsilon^{(N)}(\theta)$$

with

$$\|E_\varepsilon^{(N)}\|_\rho \leq C|\varepsilon|^{N+1}.$$  

**Remark 30.** The first step of the Newton method could start with $K^{(N_0)}$, $\mu^{(N_0)}$, given by theorem 19, for some $N_0$.

Newton’s method consists in finding corrections $\Delta_\varepsilon$, $\sigma_\varepsilon$ to $K_\varepsilon^{(N)}$ and $\mu_\varepsilon^{(N)}$ such that the linear approximation of equation $(3.1)$ associated to $K_\varepsilon^{(N)} + \Delta_\varepsilon, \mu_\varepsilon^{(N)} + \sigma_\varepsilon$ reduces the error up to quadratic terms. Taking into account that

$$f_{c,\mu;\varepsilon} \circ (K + \Delta) = f_{c,\mu} \circ K + \left[Df_{c,\mu} \circ K\right] \Delta + \left[D\mu_{c,\mu} \circ K\right] \sigma + O(\|\Delta\|^2) + O(\|\sigma\|^2)$$

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the Newton equation is
\[
\left[Df_{\varepsilon \mu[\varepsilon]} \circ K_{\varepsilon}[\varepsilon] \right] \Delta \varepsilon - \Delta \varepsilon \circ T_\nu + \left[Df_{\varepsilon \mu[\varepsilon]} \circ K_{\varepsilon}[\varepsilon] \right] \sigma \varepsilon = -E_N^\varepsilon \tag{4.5}
\]

Equation (4.5) is not easy to solve due to the fact that \(Df_{\varepsilon \mu[\varepsilon]} \circ K_{\varepsilon}[\varepsilon]\) is not constant. Following an approach similar to that in [CCdlL13], we will not solve (4.5) exactly but we will find approximate solutions that will reduce quadratically the error. The idea is to approximate the solution of (4.5) using the geometric identities introduced in section 2.6. Considering the change of variables
\[
\Delta \varepsilon = M_{\varepsilon[\varepsilon]} \varepsilon, \tag{4.6}
\]
where \(M_{\varepsilon[\varepsilon]}\) is as in (2.8) computed from \(K_{\varepsilon}[\varepsilon]\). Using (2.11) one obtains that (4.5) is equivalent to
\[
M_{\varepsilon[\varepsilon]} \circ T_\nu \left[ \begin{pmatrix} \text{Id} & S_{\varepsilon[\varepsilon]} \lambda(\varepsilon) \text{Id} \\ 0 & \lambda(\varepsilon) \text{Id} \end{pmatrix} \right] \varepsilon - W_\varepsilon \circ T_\nu + \left(Df_{\varepsilon \mu[\varepsilon]} \circ K_{\varepsilon}[\varepsilon]\right) \sigma \varepsilon = -E_N^\varepsilon - R_{\varepsilon[\varepsilon]} \varepsilon \tag{4.7}
\]
where \(R_{\varepsilon[\varepsilon]}\) is the error (2.13) and \(S_{\varepsilon[\varepsilon]}\) is given in (2.12), both computed from \(K_{\varepsilon}[\varepsilon]\). That is
\[
M_{\varepsilon[\varepsilon]} \equiv \left[Df_{\varepsilon[\varepsilon]} \mid J^{-1} \circ K_{\varepsilon} \circ J^{-1} \circ K_{\varepsilon} \right] \sim O(|\varepsilon|^0) \tag{4.8}
\]

\[
S_{\varepsilon[\varepsilon]} \equiv \left \{ P_{\varepsilon[\varepsilon]} \right \}^T Df_{\varepsilon \mu[\varepsilon]} \circ K_{\varepsilon}[\varepsilon] J^{-1} \circ K_{\varepsilon}[\varepsilon] P_{\varepsilon[\varepsilon]} - \lambda(\varepsilon) \lambda_{\varepsilon}[\varepsilon]^T \Gamma_{\varepsilon}[\varepsilon] \lambda_{\varepsilon}[\varepsilon] \sim O(|\varepsilon|^0) \tag{4.9}
\]
\[
\lambda_{\varepsilon}[\varepsilon] \equiv \left \{ \left( \left( Df_{\varepsilon[\varepsilon]} \right) \right)^T Df_{\varepsilon[\varepsilon]} \right \}^{-1} \sim O(|\varepsilon|^0), \tag{4.10}
\]
\[
P_{\varepsilon[\varepsilon]} \equiv Df_{\varepsilon}[\varepsilon] \lambda_{\varepsilon}[\varepsilon], \Gamma_{\varepsilon}[\varepsilon] \equiv Df_{\varepsilon[\varepsilon]} J^{-1} \circ K_{\varepsilon}[\varepsilon] Df_{\varepsilon[\varepsilon]} \tag{4.11}
\]

Since we expect both \(W_\varepsilon\) and \(R_{\varepsilon[\varepsilon]}\) to be estimated by \(E_N^\varepsilon\), see (5.5) and (5.15), the term \(W_\varepsilon R_{\varepsilon[\varepsilon]}\) is quadratic in \(E_N^\varepsilon\), thus, we expect that omitting this term in (4.7) will not change the quadratic nature of the method.

In order to be able to get estimates of solutions of cohomology equations of the form (4.3) instead of considering the whole error \(E_N^\varepsilon = \sum_{j=N+1}^\infty E_{\varepsilon^j}\) we only consider a truncation of this series, that is, we only consider \(E_{\varepsilon^{2N}} = \sum_{j=N+1}^{2N} E_{\varepsilon^j}\).

Taking the above into account our quasi Newton step consists in solving the following equation
\[
M_{\varepsilon[\varepsilon]} \circ T_\nu \left[ \begin{pmatrix} \text{Id} & S_{\varepsilon[\varepsilon]} \lambda(\varepsilon) \text{Id} \\ 0 & \lambda(\varepsilon) \text{Id} \end{pmatrix} \right] \varepsilon - W_\varepsilon \circ T_\nu + \left(Df_{\varepsilon \mu[\varepsilon]} \circ K_{\varepsilon}[\varepsilon]\right) \sigma \varepsilon = -E_{\varepsilon^{2N}} \tag{4.12}
\]
Remark 31. The choice of the truncation $E^{N,2N}_{\varepsilon}$ in (4.12) has two very important implications for the proof of our result. The first one is that this will yield a new approximate solution which reduces the error quadratically, as a function of $\varepsilon$. Moreover, our model example, the dissipative standard map (1.1), will satisfy hypothesis HTP1 and HTP2 in theorem B due to the fact that the truncation is made. See appendix A.

In order to construct a solution of equation (4.12), we follow a similar approach as in [CCdlL13]. Defining
\[ E^{(N,2N)}_{\varepsilon} := (M^{\varepsilon}_{N} \circ T_{\varepsilon})^{-1} E^{(N,2N)}_{\varepsilon} \sim O(||\varepsilon||^{N+1}) \]  
and writing $E^{(N,2N)}_{\varepsilon} \equiv \begin{pmatrix} E^{(N,2N)}_{\varepsilon,1} \\ E^{(N,2N)}_{\varepsilon,2} \end{pmatrix}$, where $E^{(N,2N)}_{\varepsilon,1}$ and $E^{(N,2N)}_{\varepsilon,2}$ are the first and last $d$ rows of the $2d \times 1$ matrix $E^{(N,2N)}_{\varepsilon}$. Similarly, write $A^{\varepsilon}_{N} = \begin{pmatrix} A^{\varepsilon}_{N,1} \\ A^{\varepsilon}_{N,2} \end{pmatrix}$ and $W_{\varepsilon} = \begin{pmatrix} W_{\varepsilon,1} \\ W_{\varepsilon,2} \end{pmatrix}$. Then (4.12) can be written in components as
\[ W_{\varepsilon,1} - W_{\varepsilon,1} \circ T_{\varepsilon} = -S^{(N,2N)}_{\varepsilon} W_{\varepsilon,2} - E^{(N,2N)}_{\varepsilon,1} - A^{\varepsilon}_{N,1} \sigma_{\varepsilon} \]
\[ \lambda(\varepsilon) W_{\varepsilon,2} - W_{\varepsilon,2} \circ T_{\varepsilon} = -E^{(N,2N)}_{\varepsilon,2} - A^{\varepsilon}_{N,2} \sigma_{\varepsilon}. \]

Denoting $\overline{W}_{\varepsilon,i}$ as the average of $W_{\varepsilon,i}$ with respect to $\theta$, and $(W_{\varepsilon,i})^{0} = W_{\varepsilon,i} - \overline{W}_{\varepsilon,i}$, $i = 1, 2$, we can divide the system above into two systems, one for the average and another one for the no-average part, that is
\[ 0 = -S^{(N,2N)}_{\varepsilon} \overline{W}_{\varepsilon,2} - S^{(N,2N)}_{\varepsilon} (W_{\varepsilon,2})^{0} - E^{(N,2N)}_{\varepsilon,1} - A^{\varepsilon}_{N,1} \sigma_{\varepsilon} \]
\[ (W_{\varepsilon,1})^{0} - (W_{\varepsilon,1})^{0} \circ T_{\varepsilon} = -\left( S^{(N,2N)}_{\varepsilon} W_{\varepsilon,2} \right)^{0} - \left( E^{(N,2N)}_{\varepsilon,1} \right)^{0} - \left( A^{\varepsilon}_{N,1} \right)^{0} \sigma_{\varepsilon} \]
\[ \lambda(\varepsilon) (W_{\varepsilon,2})^{0} - (W_{\varepsilon,2})^{0} \circ T_{\varepsilon} = -\left( E^{(N,2N)}_{\varepsilon,2} \right)^{0} - \left( A^{\varepsilon}_{N,2} \right)^{0} \sigma_{\varepsilon}. \]

In order to uncouple systems (4.17) and (4.18) we consider $(W_{\varepsilon,2})^{0}$ as an affine function of $\sigma_{\varepsilon}$, due to (4.18). That is,
\[ (W_{\varepsilon,2})^{0} = (B_{a,\varepsilon})^{0} + (B_{b,\varepsilon})^{0} \sigma_{\varepsilon}, \]
where $(B_{a,\varepsilon})^{0}$ and $(B_{b,\varepsilon})^{0}$ are defined as the solutions of
\[ \lambda(\varepsilon) (B_{a,\varepsilon})^{0} - (B_{a,\varepsilon})^{0} \circ T_{\varepsilon} = -\left( E^{(N,2N)}_{\varepsilon,2} \right)^{0} \]
\[ \lambda(\varepsilon) (B_{b,\varepsilon})^{0} - (B_{b,\varepsilon})^{0} \circ T_{\varepsilon} = -\left( A^{\varepsilon}_{N,2} \right)^{0}. \]
Due to HTP1, and applying lemma 29, equations (4.20) and (4.21) can be solved and we can get estimates in balls with center at $\varepsilon = 0$. Once that (4.20) and (4.21) are solved, and using (4.19), system (4.17) can be written as
\[
\left( \frac{S^{[\leq N]}_\varepsilon}{ \varepsilon^3 \text{Id} } + \frac{S^{[\leq N]}_\varepsilon (B_{\varepsilon,2}^0) + \tilde{A}^{N,1}_\varepsilon}{A^{N,2}_\varepsilon} \right) \left[ W_\varepsilon (W_\varepsilon) - \left( - \frac{S^{[\leq N]}_\varepsilon (B_{\varepsilon,2}^0) - \tilde{E}^{N,2N}_\varepsilon}{E^{N,2N}_\varepsilon} \right) \right] = - \varepsilon^3 \theta \]

(4.22)

**Remark 32.** Due to HND in theorem 19 the matrix in the left-hand side of (4.22) is invertible at \( \varepsilon = 0 \). By the continuity of the determinant, equation (4.22) can be solved for \( \varepsilon \) small enough and the inverse is analytic in \( \varepsilon \).

Thus, (4.19) and (4.22) yield \( \sigma_\varepsilon \sim O \left( |\varepsilon|^{N+1} \right) \) and \( W_{\varepsilon,2} = (W_{\varepsilon,2})^0 + \bar{W}_{\varepsilon,2} \sim O \left( |\varepsilon|^{N+1} \right) \). It remains to find \( W_{\varepsilon,1} \), this can be done by solving the equation

\[
(W_{\varepsilon,1})^0 - (W_{\varepsilon,1})^0 \circ T_\varepsilon = - \left( S^{[\leq N]}_\varepsilon W_{\varepsilon,1} \right)^0 - \left( E^{N,2N}_\varepsilon \right)^0 - \left( \tilde{A}^{N,1}_\varepsilon \right)^0 \sigma_\varepsilon,
\]

(4.23)

which can be done due to lemma 26. To fulfill the normalization condition (3.3) and obtain uniqueness of the coefficients of the perturbative expansions, \( \bar{W}_{\varepsilon,1} \) is chosen as

\[
\bar{W}_{\varepsilon,1} = - \left( \int_{\mathbb{T}^d} \left[ M_0^{-1}(\theta)DK^{[\leq N]}_\varepsilon \right]_d d\theta \right)^{-1} \times \int_{\mathbb{T}^d} \left[ M_0^{-1}(\theta) \left( DK^{[\leq N]}_\varepsilon \right)^0 \left( V^{[\leq N]}_\varepsilon W_{\varepsilon,1} \right) + V^{[\leq N]}_\varepsilon W_{\varepsilon,2} \right]_d d\theta,
\]

(4.24)

where \( V^{[\leq N]}_\varepsilon = J^{-1} \circ K^{[\leq N]}_\varepsilon DK^{[\leq N]}_\varepsilon \tilde{A}^{[\leq N]}_\varepsilon \) is the second column of the matrix \( M^{[\leq N]}_\varepsilon \), see remark 17.

**Remark 33.** If \( K^{[\leq N]}_\varepsilon \) satisfies the normalization (3.3), then the new approximation \( K^{[\leq N]}_\varepsilon + \Delta_\varepsilon \) will satisfy (3.3) if the correction satisfies

\[
\int_{\mathbb{T}^d} M_0^{-1}(\theta) \Delta_\varepsilon(\theta) d\theta = 0.
\]

Since \( \Delta_\varepsilon = M^{[\leq N]}_\varepsilon W_\varepsilon = DK^{[\leq N]}_\varepsilon W_{\varepsilon,1} + V^{[\leq N]}_\varepsilon W_{\varepsilon,2} = DK^{[\leq N]}_\varepsilon \left( (W_{\varepsilon,1})^0 + \bar{W}_{\varepsilon,1} \right) + V^{[\leq N]}_\varepsilon W_{\varepsilon,2} \), (4.24) follows from the fact that \( \int_{\mathbb{T}^d} \left[ M_0^{-1}(\theta)DK^{[\leq N]}_\varepsilon \right]_d d\theta = \int_{\mathbb{T}^d} \left[ M_0^{-1}(\theta)DK^{[\leq N]}_\varepsilon \right]_d d\theta \int_{\mathbb{T}^d} \left[ M_0^{-1}(\theta)DK^{[\leq N]}_\varepsilon \right]_d d\theta \) is invertible, for \( \varepsilon \) small enough, due to the fact that \( DK^{[\leq N]}_\varepsilon(\theta) \) is a perturbation of \( DK_0(\theta) \) and \( [M_0^{-1}(\theta)DK_0(\theta)]_d = I_d \), because \( M_0(\theta) = [DK_0(\theta) | V_0(\theta)] \).

This yields, \( W_{\varepsilon,1} = (W_{\varepsilon,1})^0 + \bar{W}_{\varepsilon,1} \sim O \left( |\varepsilon|^{N+1} \right) \) and thus

\[
\Delta_\varepsilon = M^{[\leq N]}_\varepsilon W_\varepsilon \sim O \left( |\varepsilon|^{N+1} \right) \quad \text{and} \quad \sigma_\varepsilon \sim O \left( |\varepsilon|^{N+1} \right).
\]

(4.25)

which means that \( \Delta_\varepsilon \equiv \sum_{n=N+1}^{\infty} \Delta_n \varepsilon^n \) and \( \sigma_\varepsilon \equiv \sum_{n=N+1}^{\infty} \sigma_n \varepsilon^n \). Finally, we take the corrections as

\[
\Delta^{(N,2N)}_\varepsilon \equiv \sum_{n=N+1}^{2N} \Delta_n \varepsilon^n \quad \text{and} \quad \sigma^{(N,2N)}_\varepsilon \equiv \sum_{n=N+1}^{2N} \sigma_n \varepsilon^n.
\]

(4.26)

Therefore, the new approximation is chosen as

\[
K^{[\leq N]}_\varepsilon := K^{[\leq N]}_\varepsilon + \Delta^{(N,2N)}_\varepsilon \quad \text{and} \quad \mu^{[\leq N]}_\varepsilon := \mu^{[\leq N]}_\varepsilon + \sigma^{(N,2N)}_\varepsilon.
\]

(4.27)
Remark 34. Notice that, due to lemma 29, the solutions of (4.20) and (4.21) will satisfy \((B_{\alpha\varepsilon})^0 \sim O(|\varepsilon|^N)\) and \((B_{\beta\varepsilon})^0 \sim O(|\varepsilon|^0)\), because \((\hat{E}_{\varepsilon_2}^{(N,2N)})^0 \sim O(|\varepsilon|^N)\) and \((\hat{A}_{\varepsilon_2}^{N})^0 \sim O(|\varepsilon|^0)\). Moreover, (4.22) implies that \(W_{\varepsilon_2} \sim O(|\varepsilon|^N)\) and \(\sigma \sim O(|\varepsilon|^N)\). Thus, \(W_{\varepsilon_2} \sim O(|\varepsilon|^N)\) and similarly \(W_{\varepsilon_1} \sim O(|\varepsilon|^N)\) which implies \(\Delta \sim O(|\varepsilon|^N)\).

4.3. Algorithm for the iterative step

The procedure described above leads to algorithm 35 for a given Diophantine vector \(\omega\) and assuming that we are given an analytic family \(f_{\varepsilon,\omega}\). Some steps in the algorithm are denoted as \(p \leftarrow q\), meaning that the quantity \(q\) is assigned to the variable \(p\).

Algorithm 35. Given \(K_1^{[\varepsilon,N]} : \mathbb{T}^n \to \mathcal{M}, \mu_\varepsilon^{[\varepsilon,N]} \in \mathbb{R}^d\). We perform the following computations:

1. \(E_\varepsilon^{[\varepsilon,N]} \leftarrow f_{\varepsilon,\omega}^{[\varepsilon,N]} \circ K_1^{[\varepsilon,N]} - K_1^{[\varepsilon,N]} \circ \omega\)
2. \(E_\varepsilon^{(N,2N)}\) obtained from \(E_\varepsilon^{[\varepsilon,N]}\) by truncation
3. \(\alpha_\varepsilon \leftarrow DK_1^{[\varepsilon,N]}\)
4. \(N_\varepsilon \leftarrow [\alpha_\varepsilon J]^{-1}\)
5. \(V_\varepsilon \leftarrow J^{-1} \circ K_1^{[\varepsilon,N]} \circ \alpha_\varepsilon \circ N_\varepsilon\)
6. \(M_\varepsilon \leftarrow [\alpha_\varepsilon | V_\varepsilon]\)
7. \(\beta_\varepsilon \leftarrow (M_\varepsilon \circ \omega)^{-1}\)
8. \(E_\varepsilon^{(N,2N)} \leftarrow \beta_\varepsilon E_\varepsilon^{(N,2N)}\)
9. \(P_\varepsilon \leftarrow \alpha_\varepsilon N_\varepsilon\)

\[
\Gamma_\varepsilon \leftarrow \alpha_\varepsilon J^{-1} \circ K_1^{[\varepsilon,N]} \circ \alpha_\varepsilon
\]

\[
S_\varepsilon \leftarrow (P_\varepsilon \circ \omega)^{-1/2} Df_{\mu_\varepsilon^{[\varepsilon,N]} \circ K_1^{[\varepsilon,N]} \circ \omega} J^{-1} \circ K_1^{[\varepsilon,N]} \circ P_\varepsilon
\]

\[
\lambda(\varepsilon)(N_\varepsilon \circ \omega)^{-1} \Gamma_\varepsilon \circ T_\omega (N_\varepsilon \circ T_\omega)
\]

\[
\hat{A}_\varepsilon \leftarrow \beta_\varepsilon \mu f_{\mu_\varepsilon^{[\varepsilon,N]} \circ K_1^{[\varepsilon,N]} \circ \omega}
\]

10. \((B_{\alpha\varepsilon})^0 \) solves \(\lambda(\varepsilon)(B_{\alpha\varepsilon})^0 - (B_{\alpha\varepsilon})^0 \circ \omega = -\left(E_\varepsilon^{(N,2N)}\right)^0\)

\((B_{\beta\varepsilon})^0 \) solves \(\lambda(\varepsilon)(B_{\beta\varepsilon})^0 - (B_{\beta\varepsilon})^0 \circ \omega = -\left(\hat{A}_\varepsilon\right)^0\)

11. Find \(W_{\varepsilon_2}, \sigma\) by solving

\[
\begin{pmatrix}
S_\varepsilon & S_\varepsilon(B_{\beta\varepsilon})^0 + A_{\varepsilon_2}
\end{pmatrix}
\begin{pmatrix}
W_{\varepsilon_2}
\end{pmatrix}
= \begin{pmatrix}
-S_\varepsilon(B_{\alpha\varepsilon})^0 - E_\varepsilon^{(N,2N)}
\end{pmatrix}
\]

12. \((W_{\varepsilon_2})^0 = (B_{\alpha\varepsilon})^0 + (B_{\beta\varepsilon})^0 \sigma_\varepsilon\)

13. \(W_{\varepsilon_2} = (W_{\varepsilon_2})^0 + W_{\varepsilon_2} \sim O(|\varepsilon|^N)\)
\[(14) \,(W_{\varepsilon,1})^0 + (S_{\varepsilon} W_{\varepsilon,2})^0 - (W_{\varepsilon,1})^0 \circ T_{\varepsilon} = (E_{\alpha,2}^{N,2N})^0 - (\bar{A}_{-1})^0\]

\[(15) \, W_{\varepsilon,1} = -\left(\int_{\mathbb{R}^d} [M_0^{-1} \alpha_{1}]_1 d\theta \right)^{-1} \int_{\mathbb{R}^d} [M_0^{-1} (\alpha_{1} (W_{\varepsilon,1})^0 + V_{\varepsilon} W_{\varepsilon,2})]_1 d\theta\]

\[(16) \, W_{\varepsilon,1} = (W_{\varepsilon,1})^0 + \bar{W}_{\varepsilon,1} \sim O\left(|\varepsilon|^{N+1}\right)\]

\[(17) \, \Delta_{\varepsilon} \leftarrow M_{\varepsilon} W_{\varepsilon}\]

\[(18) \, K_{1}^{[\leq 2N]} \leftarrow K_{1}^{[\leq N]} + \Delta_{1}^{(N,2N)}\]

It is worth noting that all the operations in algorithm 35 could be implemented in a few lines in a high level computer language.

**Remark 36.** Note that algorithm 35 involves only algebraic operations, compositions, derivatives, truncations, and solving cohomology equations. This implies that if we start with analytic functions then the output will be an analytic function.

**Remark 37.** Note that at each step of the iterative procedure obtained by the quasi Newton method the input will be polynomials of degree \(N\) in \(\varepsilon\), \(K_{1}^{[\leq N]} \equiv \sum_{n=0}^{N} K_{n} e_{n}\), and \(\mu_{1}^{[\leq N]} = \sum_{n=0}^{N} \mu_{n} e_{n}\). The output will be polynomials of degree \(2N\) in \(\varepsilon\) given by

\[K_{2}^{[\leq 2N]} : K_{1}^{[\leq N]} + \Delta_{2}^{(N,2N)}\] and \(\mu_{2}^{[\leq 2N]} : = \mu_{1}^{[\leq N]} + \sigma_{2}^{(N,2N)}\).

Since, by construction, \(\Delta_{2}^{(N,2N)} \sim O\left(|\varepsilon|^{N+1}\right)\) and \(\sigma_{2}^{(N,2N)} \sim O\left(|\varepsilon|^{N+1}\right)\), the first \(N\) coefficients \(K_{1}, K_{2}, \ldots, K_{N}\) of the expansion of \(K_{1}^{[\leq 2N]}\) will be the same coefficients of \(K_{1}^{[\leq N]}\) and they will not change for any of the next steps. The same also happens for the coefficients of \(\mu_{1}^{[\leq 2N]}\). This is a crucial step for proving theorem B, since due to the fact that the coefficient up to order \(N\) do not change after \(\log_{2}(N)\) steps of the modified Newton method, one can use Cauchy estimates in the domains given by lemma 29 after \(\log_{2}(N)\) steps to obtain estimates on the \(N\)th coefficient.

**Remark 38.** To iterate the modified Newton method in algorithm 35 it is needed that the new error \(E_{2N}^{2N}\) obtained using the new approximations \(K_{1}^{[\leq 2N]} = K_{1}^{[\leq N]} + \Delta_{2}^{(N,2N)}\) and \(\mu_{1}^{[\leq 2N]} = \mu_{1}^{[\leq N]} + \sigma_{2}^{(N,2N)}\) satisfies \(E_{2N}^{2N} \sim O\left(|\varepsilon|^{2N+1}\right)\). This is a consequence of the fact that the new error is quadratic in the original error, as an expansion on \(\varepsilon\), and this is verified in proposition 47.

### 5. Estimates for the iterative step

In this section we present the estimates for the corrections given by the Newton step described in section 4, these estimates are obtained by following the steps in algorithm 35. Throughout this section we consider maps in the spaces \(A_{\mu,\gamma}\). In the following we will be dealing with equations of the form (4.3) which, according to lemma 29, can be solved if

\[|\varepsilon| \leq \gamma_{N} := \left(\frac{\nu}{2}\right)^{1/\alpha} \frac{1}{(aN)^{1/\alpha}}.\]
where $aN$ is the degree of the trigonometric polynomial in the right-hand side of (4.3).

5.1. Estimate for the reducibility error

The following lemma provides an estimate for the error in the approximate reducibility given by $R_{\leq N}$ as in (2.13) computed from $K_{\leq N}$. The estimates are obtained by studying qualitatively the geometric identities introduced in section 2.6 and taking into account the uniformity on the variable $\varepsilon$.

**Lemma 39.** Let $N \in \mathbb{N}$, $\omega \in \mathcal{D}(\nu, \tau)$ and $f_{\varepsilon,M} : \mathcal{M} \to \mathcal{M}$ be a family of analytic conformally symplectic maps, with $f_{\varepsilon,M}^* \Omega = \lambda(\varepsilon)\Omega$, $\mu \subseteq \mathbb{C}$, $d$. Let $K_{\leq N} \in \mathcal{A}_{\gamma_N}$ such that $K_{\leq N} : \mathbb{T}^d \to \mathcal{M}$ is an embedding for any $|\varepsilon| \leq \gamma_N$. Assume also that, for any $|\varepsilon| \leq \gamma_N$.

(a) $K_{\leq N}(\mathbb{T}^d) \subset \text{Domain}(f_{\varepsilon,M|\leq N})$ and that there exist $\xi \geq 0$ such that

$$\text{dist} \left( K_{\leq N}(\mathbb{T}^d), \partial \text{Domain}(f_{\varepsilon,M|\leq N}) \right) \geq \xi > 0$$

$$\text{dist} \left( \mu_{\varepsilon,M|\leq N}, \partial \Lambda \right) \geq \xi > 0.$$

(b) The approximate invariance equation holds

$$f_{\varepsilon,M|\leq N} \circ K_{\leq N} - K_{\leq N \circ T_\varepsilon} \sim \mathcal{O}(|\varepsilon|^{N+1})$$

(c) $\nu^{-1}(aN)^{-d} \left( |E^N| \right)_{\rho/\gamma_N} \ll 1$ (5.2)

(d) $\text{H} \text{T} \text{P} 2$ The $d \times d$ matrix

$$E_{\varepsilon,M}^N(\theta) \equiv DK_{\varepsilon,G}^N(\theta + \omega) \mathcal{J} \circ K_{\varepsilon,M|\leq N}(\theta + \omega)$$

$$- D(f_{\varepsilon,M|\leq N} \circ K_{\varepsilon,M|\leq N}(\theta)) \mathcal{J} \circ (f_{\varepsilon,M|\leq N} \circ K_{\varepsilon,M|\leq N}(\theta)) \mathcal{J} \circ (f_{\varepsilon,M|\leq N} \circ K_{\varepsilon,M|\leq N}(\theta))$$

$$\times D(f_{\varepsilon,M|\leq N} \circ K_{\varepsilon,M|\leq N}(\theta))$$

is a trigonometric polynomial of degree less than $aN$.

Then

$$R_{\leq N} \sim \mathcal{O}(|\varepsilon|^{N+1})$$

and for any $0 < \delta \leq \rho$ we have

$$\left\| R_{\leq N} \right\|_{\rho, \gamma_N} \leq C \nu^{-1}(aN)^{-d} \left\| E^N \right\|_{\rho, \gamma_N}.$$ (5.5)

where $C = C(d, \left\| DK_{\leq N} \right\|_{\rho, \gamma_N}, \left\| \Lambda \right\|_{\rho, \gamma_N}, \left\| J \circ K_{\leq N} \right\|_{\rho, \gamma_N}).$

**Proof.** Writing $R_{\varepsilon,M|\leq N}(\theta)$ in terms of $K_{\varepsilon,M|\leq N}$ as in (2.13) yields

$$R_{\varepsilon,M|\leq N}(\theta) = \left[ DE_{\varepsilon,M}^N(\theta) | \Lambda_{\varepsilon,M|\leq N}(\theta + \omega)(B(\theta) - \lambda(\varepsilon)\text{id}) \right]$$

$$+ DK_{\varepsilon,M|\leq N}(\theta + \omega) \left( \tilde{S}_\varepsilon(\theta) - S_{\varepsilon,M|\leq N}(\theta) \right)$$
Then, by lemma 29 and with the matrix associated to the symplectic form, see section 2. It is easy to estimate the first column

\[ K(c,\lambda) \]  

significantly simpler than in the general case.

Note that when the matrix \( R_{\mathcal{L}} \) which yields

\[ \tilde{S}_{\mathcal{L}}(\theta) = S_{\mathcal{L}}(\theta) + (\theta + \omega)(B_{\mathcal{L}}(\theta) - \lambda(\varepsilon) \text{Id}) \]

is the pull back \((K(\varepsilon)^{\text{sym}})^{\mathcal{L}}\) written in coordinates and \([\varepsilon]^{\mathcal{L}}\) as in (4.11). We recall that \( J \) is the matrix corresponding to the symplectic form, see section 2. It is easy to estimate the first column of \( R_{\mathcal{L}}^{(\varepsilon)} \) using Cauchy estimates, that is

\[ \left\| D E_{\mathcal{L},\varepsilon}^{N} (\theta) \right\|_{\rho - \delta, \gamma N} \leq C \delta^{-1} \left\| E_{\mathcal{L},\varepsilon}^{N} \right\|_{\rho} \]

As it is pointed out in remark 13, (5.7) and (5.8) define a system of equations, for \( B_{\mathcal{L}} \) and \( \tilde{S}_{\mathcal{L}} \), which can be solved as long as \( E_{\mathcal{L},\varepsilon}^{N} \) is small enough. Thus, to obtain estimates for the second column of \( R_{\mathcal{L}}^{(\varepsilon)} \) it is enough to get estimates of \( E_{\mathcal{L},\varepsilon}^{N} \). The estimate for \( E_{\mathcal{L}}^{N} \) is obtained using that \( f_{\varepsilon,\gamma}^{\mathcal{L}} \Omega = \lambda(\varepsilon) \Omega \). Note that \( E_{\mathcal{L},\varepsilon}^{N} = (K_{\varepsilon}^{(\varepsilon)} \circ T_{\varepsilon})^{\mathcal{L}} \Omega - (f_{\varepsilon,\gamma}^{\mathcal{L}} \circ K_{\varepsilon}^{(\varepsilon)})^{\mathcal{L}} \Omega \) in coordinates and, since \((f_{\varepsilon,\gamma}^{\mathcal{L}} \circ K_{\varepsilon}^{(\varepsilon)})^{\mathcal{L}} \Omega = \lambda(K_{\varepsilon}^{(\varepsilon)})^{\mathcal{L}} \Omega\), we have that \( E_{\mathcal{L}}^{N} \) satisfies the equality

\[ E_{\mathcal{L},\varepsilon}^{N} \circ T_{\omega} - \lambda(\varepsilon) E_{\mathcal{L},\varepsilon}^{N} = E_{\mathcal{L},\varepsilon}^{N} \]

Then, by lemma 29 and HTP2 we obtain

\[ \left\| E_{\mathcal{L}}^{N} \right\|_{\rho - \delta, \gamma N} \leq C\nu^{-1}(\alpha N)^{\delta - \frac{d}{2}} \left\| E_{\mathcal{L}}^{N} \right\|_{\rho - \delta, \gamma N} \]

To get estimates for \( E_{\mathcal{L}}^{N} \), we follow [CCdlL13]. If \( h \) and \( g \) are smooth maps with range in \( \mathcal{M} \), the matrix corresponding to \( h^{\top} \Omega - g^{\top} \Omega \) is

\[ Dh^{\top} J \circ h Dh - Dg^{\top} J \circ g Dh = (Dh^{\top} - Dg^{\top})J \circ h Dh - Dg^{\top} (J \circ h - J \circ g) Dh \]

\[ + Dg^{\top} J \circ g (Dh - Dg) \]

Using this formula with \( g = f_{\varepsilon,\mu_{\varepsilon}}^{\mathcal{L}} \circ K_{\varepsilon}^{(\varepsilon)}, h = K_{\varepsilon}^{(\varepsilon)} \circ T_{\omega} \) and Cauchy estimates one obtains

\[ \left\| E_{\mathcal{L},\varepsilon}^{N} \right\|_{\rho - \delta, \gamma N} \leq C \delta^{-1} \left\| E_{\mathcal{L}}^{N} \right\|_{\rho} \]

which yields \( E_{\mathcal{L},\varepsilon}^{N}, E_{\mathcal{L}}^{N} \sim \mathcal{O}(\varepsilon^{N+1}) \) and, then, \( R_{\mathcal{L}}^{(\varepsilon)} \sim \mathcal{O}(\varepsilon^{N+1}) \) and

\[ \left\| R_{\mathcal{L}}^{(\varepsilon)} \right\|_{\rho - \delta, \gamma N} \leq C \nu^{-1}(\alpha N)^{\delta - \frac{d+1}{2}} \left\| E_{\mathcal{L}}^{N} \right\|_{\rho \gamma N} \]

Note that when the matrix \( J \) is constant both HTP2 and the computations above are significantly simpler than in the general case.
Remark 40. We emphasize that, if \( K_0 \) satisfies \( K_0 \circ T_\omega = f_{0,\mu_0} \circ K_0 = 0 \) then \( \Delta K_\omega(\theta)^{-1} J \circ K_0 \Delta K_\omega(\theta) = 0 \) and \( K_0(\mathcal{T}^N) \) is a Lagrangian manifold, see \([CCdlL13]\). This implies that the spaces \( \text{range} (\Delta K_\omega(\theta)) \) and \( \text{range} (J^{-1} \circ K_0(\theta) \Delta K(\theta)) \) are transversal and this condition makes \( M_0(\theta) \) a linear isomorphism. Note that if \( E^N_{\rho} \) in (5.9) represents the error of the Lagrangian character of \( K^\varepsilon \), then, if \( E^N_\rho \) is small enough the spaces \( \text{Range}(\Delta K^\varepsilon(\theta)) \) and \( \text{Range}(J^{-1} \circ K^\varepsilon(\theta) \Delta K^\varepsilon(\theta)) \) will be transversal and the matrix \( M^\varepsilon_{\rho} \) will define a linear isomorphism. This transversality will be obtained if (5.2) is satisfied and it is given by (5.11) and (5.12).

5.2. Estimates for the corrections

In this section we obtain estimates for the corrections \( \Delta^{(N,2N)} \) and \( \sigma^{(N,2N)} \), this estimates are obtained by following the steps in algorithm 35. First, lemma 41, we obtain estimates for the corrections \( \Delta \), \( \sigma \) and then, using Cauchy estimates, we obtain estimates for the truncations \( \Delta^{(N,2N)} \), \( \sigma^{(N,2N)} \), corollary 42.

Consider \( C \subseteq \mathbb{C}^d / \mathbb{Z}^d \times \mathbb{C}^d \) the complexification of \( \mathcal{M} = \mathbb{T}^d \times B \).

Lemma 41. Let \( \varepsilon \in \mathbb{N}_0 \), \( 0 < \rho < 1 \), and \( \delta \) such that \( 0 < 2\delta < \rho \). Assume that for any \( \varepsilon \in C \), such that \( |\varepsilon| < \gamma \), \( f_{\varepsilon,\mu}^{(N,\varepsilon)} : C \to C \) is an analytic conformally symplectic map with \( f_{\varepsilon,\mu}^{(N,\varepsilon)} \Omega = \lambda(\varepsilon) \Omega \). Assume also that \( K^{(N,\varepsilon)} \in \mathcal{A}_{\rho,\gamma} \) is such that \( K^{(N,\varepsilon)} : T_{\rho} \to \mathbb{C}^d / \mathbb{Z}^d \times \mathbb{C}^d \) is an embedding. Assume also that for any \( |\varepsilon| < \gamma_N \) we have the following:

(a) \( K^{(N,\varepsilon)}(T_{\rho}) \subset \text{Domain}(f_{\varepsilon,\mu}^{(N,\varepsilon)}) \) and that there exist \( \xi \geq 0 \) such that

\[
\text{dist} \left( K^{(N,\varepsilon)}(T_{\rho}), \partial \text{Domain}(f_{\varepsilon,\mu}^{(N,\varepsilon)}) \right) \geq \xi > 0
\]

(b) \( \text{HND} \). The following non-degeneracy condition holds:

\[
\det \left( S^{(N,\varepsilon)}_{\varepsilon,\mu} + \tilde{A}^{(N,\varepsilon)}_{\varepsilon,\mu} \right) \neq 0
\]

(c) For any \( N \in \mathbb{N} \), the matrices \( \left( \tilde{E}^{(N,2N)}_{\varepsilon,\mu} \right)^0 \) and \( \left( \tilde{A}^{(N,2N)}_{\varepsilon,\mu} \right)^0 \) defined in (4.13) and (4.14), are trigonometric polynomials of degree less or equal than \( aN \).

Then, for any \( 0 < r < 1 \) we have

\[
W_\varepsilon \sim O \left( |\varepsilon|^{N+1} \right), \quad \sigma_\varepsilon \sim O \left( |\varepsilon|^{N+1} \right)
\]

\[
\|W\|_{r^{-\delta},\rho,\gamma} \leq C \nu^{-3} (aN)^2 \delta^{-2-(r+3d)} \frac{r^{N+1}}{1-r} \|E^N\|_{\rho,\gamma}
\]

and

\[
\sup_{|\varepsilon| \leq \rho} |\sigma_\varepsilon| \leq C \nu^{-1} (aN)^2 \delta^{-2} \frac{r^{N+1}}{1-r} \|E^N\|_{\rho,\gamma}
\]

where \( C = (d, \|DK^{(N,\varepsilon)}\|_{\rho,\gamma}, \|M^{(N,\varepsilon)}\|_{\rho,\gamma}^{-1}, \|M^{(N,\varepsilon)}\|_{\rho,\gamma}, \|\Lambda^{(N,\varepsilon)}\|_{\rho,\gamma}, \mathcal{T}^N) \) and \( \mathcal{T}^N \) is defined in (5.20).
Proof. Given that \((E^{(N,2N)})_0^0\) and \((\tilde{A}_{N,2})^0\) are trigonometric polynomials, by lemma 29, (4.20), and (4.21); \(B_a\) and \(B_0\) satisfy the following estimates
\[
\|B_a\|_{\rho-\delta, r_N} \leq C\nu^{-1}(aN)^\gamma \delta^{-d} \left\|E_2^{(N,2N)}\right\|_{\rho, r_N}
\]
and similarly
\[
\|B_b\|_{\rho-\delta, r_N} \leq C\nu^{-1}(aN)^\gamma \delta^{-d} \left\|A^N\right\|_{\rho, r_N}.
\]
Taking into account that \(W_2 = (W_2)^0 + \tilde{W}_2^0\) and \((W_2) = (B_2) + \sigma(B_0)\), to have estimates for \(W_2\) we need estimates for \(\tilde{W}_2\) and \(\sigma\). Now, according to (4.22) we have
\[
\begin{pmatrix}
W_{2,2} \\
\sigma_{\epsilon}
\end{pmatrix}
= \begin{pmatrix}
S_{\epsilon^{\leq N}} \frac{S_{\epsilon^{\leq N}}(B_{\alpha_{\epsilon}})^0 + A_{\epsilon,1}^N}{A_{\epsilon,2}^N} \\
\epsilon^3 \text{Id}
\end{pmatrix}^{-1}
\times \begin{pmatrix}
-S_{\epsilon^{\leq N}} \frac{B_{\alpha_{\epsilon}})^0 - E_{\epsilon,2}^{2N}}{E_{\epsilon,1}^{2N}}
\end{pmatrix},
\]

denoting
\[
T_{\epsilon}^N := \left\|\begin{pmatrix}
S_{\epsilon^{\leq N}} \frac{S_{\epsilon^{\leq N}}(B_{\alpha_{\epsilon}})^0 + A_{\epsilon,1}^N}{A_{\epsilon,2}^N} \\
\epsilon^3 \text{Id}
\end{pmatrix}^{-1}\right\|
\]
and
\[
T_\epsilon^N = \sup_{|\epsilon| \leq r_N} T_{\epsilon}^N
\]
from (5.19) we have
\[
|\sigma_{\epsilon}|_{W_{\epsilon,2}} \leq T_{\epsilon}^N \left(\left|S_{\epsilon^{\leq N}}(B_{\alpha_{\epsilon}})^0 + E_{\epsilon,1}^{2N}\right| + \left|E_{\epsilon,2}^{2N}\right|\right) \sim O(\epsilon|N+1)
\]
which yields \(\sigma_{\epsilon} \sim O(\epsilon|N+1)\) and \(\tilde{W}_{\epsilon,2} \sim O(\epsilon|N+1)\) because \((B_{\alpha_{\epsilon}})^0 \sim O(\epsilon|N+1)\) and \(E_{\epsilon,2}^{2N} \sim O(\epsilon|N+1)\).

Thus
\[
|\sigma_{\epsilon}|_{W_{\epsilon,2}} \leq T_{\epsilon}^N \left(\left|S_{\epsilon^{\leq N}}(B_{\alpha_{\epsilon}})^0 + E_{\epsilon,1}^{2N}\right| + \left|E_{\epsilon,2}^{2N}\right|\right)
\]
\[
\leq C T_{\epsilon}^N \left(\left|S_{\epsilon^{\leq N}}\right|_\rho \left\|B_{\alpha_{\epsilon}}\right\|^0_\rho + \left\|E_{\epsilon,1}^{2N}\right\|_\rho + \left\|E_{\epsilon,2}^{2N}\right\|_\rho\right)
\]
for any \(0 < \delta < \rho\). Thus, using (4.13) and (5.17) we obtain
\[
\sup_{|\epsilon| \leq r_N} |\sigma_{\epsilon}|_{W_{\epsilon,2}} \leq C\nu^{-1}(aN)^\gamma \delta^{-d} \left\|E_{\epsilon,2}^{2N}\right\|_{\rho, r_N}
\]
\[
\sup_{|\epsilon| \leq r_N} |\sigma_{\epsilon}| \leq C\nu^{-1}(aN)^\gamma \delta^{-d} \left\|E_{\epsilon,2}^{2N}\right\|_{\rho, r_N}.
\]
For \((W_2)^0 = (B_2)^0 + \sigma(B_0)^0\) we have
Thus, combining (5.22) and (5.24) we get
\[ \|W_2\|_{\rho-\delta,r_N} \leq C \nu^{-1} (aN)^{\gamma-\delta-d} \|E^{(N,2N)}\|_{\rho,r_N}. \] (5.25)

The estimates for \((W_1)^0\) come from (4.23) and lemma 26, i.e.,
\[ \|W_1\|_{\rho-2\delta,r_N} \leq C \nu^{-1} \delta^{-(\gamma+d)} \left[ \|\mathcal{S}\|_{\rho-\delta,r_N} \|W_2\|_{\rho-\delta,r_N} \right. \]
\[ + \left. \|\tilde{E}^{(N,2N)}\|_{\rho-\delta,r_N} + \sup_{|\sigma| \leq r_N} |\sigma| \|\tilde{A}^{N}\|_{\rho-\delta,r_N} \right] \]
\[ \leq C \nu^{-1} \delta^{-(\gamma+d)} \left[ \|\mathcal{S}\|_{\rho-\delta,r_N} \nu^{-2} (aN)^{\gamma-\delta-d} \|E^{(N,2N)}\|_{\rho,r_N} \right. \]
\[ + \left. \|\mathcal{M}^{(E,N)}\|^{-1}_{\rho,r_N} \|E^{(N,2N)}\|_{\rho,r_N} + \|\tilde{A}^{N}\|_{\rho,r_N} \nu^{-1} (aN)^{\gamma-\delta-d} \|E^{(N,2N)}\|_{\rho,r_N} \right] \]

that is,
\[ \|W_1\|_{\rho-2\delta,r_N} \leq C \nu^{-3} (aN)^{2\gamma-\delta-(\gamma+d)} \|E^{(N,2N)}\|_{\rho,r_N}. \] (5.26)

Finally, the estimate for \(\overline{W}_1\) comes from (4.24), that is
\[ \sup_{|\sigma| \leq r_N} \|\overline{W}_1\|_{\rho,r_N} \leq C \left( \|W_1\|_{\rho-\delta,r_N} + \|W_2\|_{\rho-\delta,r_N} \right) \]
\[ \leq C \nu^{-3} (aN)^{\gamma-\delta-(\gamma+3d)} \|E^{(N,2N)}\|_{\rho,r_N}. \] (5.27)

Putting together (5.25)–(5.27), and using the Cauchy estimates in corollary 6 yields the claimed estimate for \(W\). \(\square\)

**Corollary 42.** Assume the hypothesis of lemmas 39 and 41, for any \(0 < \delta < \rho\) and \(0 < r < 1\) we have
\[ \|\Delta^{(N,2N)}\|_{\rho-\delta,r_N} \leq C \nu^{-3} (aN)^{\gamma-\delta-(\gamma+3d)} \frac{r^{N+1}}{1-r^{1/2}} \|E^{N}\|_{\rho,r_N}. \] (5.28)
\[ \sup_{|\sigma| \leq r_N} |\sigma^{(N,2N)}_x| \leq C \nu^{-1} (aN)^{\gamma-\delta-d} \frac{r^{N+1}}{1-r^{1/2}} \|E^{N}\|_{\rho,r_N}. \] (5.29)
Moreover,\[ \|\Delta^{2N}\|_{\rho,\Delta_{2N}} \leq C\nu^{-3}(aN)^{2r} \delta^{-4(r+3d)} \frac{\nu^{N+1}}{1-r} \|E^N\|_{\rho,\gamma}, \quad (5.30) \]
\[ \sup_{|\varepsilon| \leq c_{\gamma}^{2N}} |\sigma^{2N,\varepsilon}| \leq C\nu^{-1}(aN)^{r} \delta^{-d} \frac{\nu^{N+1}}{1-r} \|E^N\|_{\rho,\gamma}, \quad (5.31) \]

**Proof.** Using the Cauchy estimates as in corollary 6 and the estimates in lemma 41 one obtains
\[ \|\Delta^{2N}\|_{\rho,\Delta_{2N}} \leq C\nu^{-3}(aN)^{2r} \delta^{-4(r+3d)} \frac{\nu^{N+1}}{1-r} \|E^N\|_{\rho,\gamma}, \]
\[ = C\nu^{-1}(aN)^{r} \delta^{-d} \frac{\nu^{N+1}}{1-r} \|E^N\|_{\rho,\gamma} \]
and
\[ \sup_{|\varepsilon| \leq c_{\gamma}^{2N}} |\sigma^{2N,\varepsilon}| \leq C\nu^{-1}(aN)^{r} \delta^{-d} \frac{\nu^{N+1}}{1-r} \|E^N\|_{\rho,\gamma}, \]
\[ = C\nu^{-1}(aN)^{r} \delta^{-d} \frac{\nu^{N+1}}{1-r} \|E^N\|_{\rho,\gamma}. \]
The other estimates are obtained similarly. \[ \square \]

5.3. *Nonlinear estimates for the quasi-Newton method*

The quasi-Newton procedure in algorithm 35 can also be described using a convenient operator notation. Defining the error functional
\[ \mathcal{E}[K, \mu] = f_{\mu} \circ K - K \circ T \omega \]
and assuming \( \Delta \) and \( \sigma \) are small enough, the Taylor expansion of \( \mathcal{E}[K + \Delta, \mu + \sigma] \) is given by
\[ \mathcal{E}[K + \Delta, \mu + \sigma] = \mathcal{E}[K, \mu] + D_1\mathcal{E}[K, \mu]\Delta + D_2\mathcal{E}[K, \mu]\sigma + \mathcal{R}[\Delta, \sigma; K, \mu], \quad (5.33) \]
where the Frechet derivatives are given by
\[ D_1\mathcal{E}[K, \mu] \Delta = (Df_{\mu} \circ K) \Delta - \Delta \circ T \omega \]
\[ D_2\mathcal{E}[K, \mu] \sigma = (Df_{\mu} \circ K) \sigma \]
and \( \mathcal{R} \) is the remainder of the Taylor expansion. Note that \( \mathcal{E}[K^{(\leq N)}, \mu^{(\leq N)}] = E^N, \) with this notation the classic Newton method would consist in finding a correction \( (\Delta_{2N}^{(k)}, \mu_{2N}^{(k)}) \) such that
\[ \mathcal{E}[K_{2N}^{(k)}, \mu_{2N}^{(k)}] + D_1\mathcal{E}[K_{2N}^{(k)}, \mu_{2N}^{(k)}] \Delta_{2N}^{(k)} + D_2\mathcal{E}[K_{2N}^{(k)}, \mu_{2N}^{(k)}] \sigma_{2N}^{(k)} = 0. \]
As it was explained before, in section 4.2, the corrections we construct with algorithm 35 do not satisfy (5.36) but they solve an approximate equation (4.12). The following lemmas give estimates for the error functional evaluated in the corrected unknowns. First, lemma 45, we give estimates for the error function evaluated in the corrected unknowns, \( \E[\hat{K}^{(N)} + \Delta, \mu^{(N)} + \sigma] \) and then, using Cauchy estimates, we obtain the estimates for the error evaluated in the truncated corrections, \( \E[\hat{K}^{(N)} + \Delta^{(N,2N)}, \mu^{(N)} + \sigma^{(N,2N)}] \), proposition 47.

**Remark 43.** We emphasize that to be able to compute \( \E[K + \Delta, \mu + \sigma] \) we need both \( \Delta \) and \( \sigma \) to be small enough, so the compositions in (5.32) are well defined. In particular \( \Delta \) and \( \sigma \) need to satisfy \( ||\Delta||, ||\sigma|| \leq \xi \) and we need to choose the domain loss. In section 6, lemma 51, we give smallness conditions on the initial error which will guarantee that the compositions will be defined at any step of the iteration. This is very standard in KAM theory.

**Lemma 44.** Assume \( 0 < r < 1 \) and \( 0 < \delta \leq \rho \). Then, under the hypothesis of lemmas 39 and 41 one has

\[
\E[\hat{K}^{(N)}_x, \mu^{(N)}_x] + D_1 \E[\hat{K}^{(N)}_x, \mu^{(N)}_x] \Delta_x + D_2 \E[\hat{K}^{(N)}_x, \mu^{(N)}_x] \sigma_x \sim O(\varepsilon^{2N+1})
\]  

(5.37) and

\[
\left\| \E[\hat{K}^{(N)}_x, \mu^{(N)}_x] + D_1 \E[\hat{K}^{(N)}_x, \mu^{(N)}_x] \Delta + D_2 \E[\hat{K}^{(N)}_x, \mu^{(N)}_x] \sigma \right\|_{r-\delta,rN} \leq \frac{r^{2N+1}}{1-r} \|E_N\|_{rN}^2 + C\nu^{-4}(a\nu)^{3\gamma} \delta^{-(\gamma+4d+1)} r^{N+1} \|E_N\|_{rN}^2.
\]  

(5.38)

**Proof.** Note that with the operator notation introduced at the beginning of this section we have \( \E(K^{(N)}_x, \mu^{(N)}_x) = E^N \). Using (2.11) and taking into account that \( \Delta_x = M^{(N)}_x W_x \) and that \( W_x \) satisfies (4.12) we have

\[
\E[\hat{K}^{(N)}_x, \mu^{(N)}_x] + D_1 \E[\hat{K}^{(N)}_x, \mu^{(N)}_x] \Delta_x + D_2 \E[\hat{K}^{(N)}_x, \mu^{(N)}_x] \sigma_x =
\]

\[
= E^N_x + (D_{f_{\rho,\mu}} \circ K^{(N)}_x) \Delta_x - \Delta_x \circ T_\omega + (D_{\rho f_{\mu}} \circ K^{(N)}_x) \sigma_x - R^{(N)}_x (M^{(N)}_x)^{-1} \Delta_x + R^{(N)}_x (M^{(N)}_x)^{-1} \Delta_x
\]

\[
= E^N_x + M^{(N)}_x \circ T_\omega \begin{pmatrix} Id & S^{(N)}_x \\ 0 & \lambda(\varepsilon) Id \end{pmatrix} (M^{(N)}_x)^{-1} \Delta_x - \Delta_x \circ T_\omega + (D_{f_{\rho,\mu}} \circ K^{(N)}_x) \sigma_x + R^{(N)}_x (M^{(N)}_x)^{-1} \Delta_x
\]

\[
= E^N_x - E^{N,2N}_x + R^{(N)}_x W_x
\]

\[
= E^{2N,\infty}_x + R^{(N)}_x W_x \sim O(\varepsilon^{2N+1}),
\]  

(5.39)

where \( E^{2N,\infty}_x = \sum_{n=2N}^{\infty} E_n \varepsilon^n \). Note that the order of \( \varepsilon \) in the last line follows from the definition of \( E^{2N,\infty}_x \), (5.4), and (5.14).

Then, using the Cauchy estimates of corollary 6, lemmas 39 and 41 one obtains
Lemma 45. Assume $0 < r < 1$ and $0 < \delta \leq \rho$. Then, under the hypothesis of lemmas 41 and 39 we have

\[
\mathcal{E}(K^\varepsilon) + \Delta_r \varepsilon [\mathcal{K}^\varepsilon + \sigma_r] \sim O \left( |\varepsilon|^{2N+1} \right) \tag{5.40}
\]

and

\[
\left\| \mathcal{E}(K^\varepsilon) + \Delta_r \varepsilon [\mathcal{K}^\varepsilon + \sigma_r] \right\|_{\rho-\delta, r\gamma_N} \leq \frac{r^{2N+1}}{1-r} \left\| E^N \right\|_{\rho, \gamma_N} + C \nu^{-6}(aN)^{3\gamma} \delta^{-(2r+6d)}
\]

where $C = C \left( \left\| D\mathcal{K}^\varepsilon \right\|_{\rho, \gamma_N}, \left\| D^2 f^\varepsilon \circ \mathcal{K}^\varepsilon \right\|_{\rho, \gamma_N}, \left\| D^2 f^\varepsilon \circ \mathcal{K}^\varepsilon \right\|_{\rho, \gamma_N} \right)$.

Proof. Note that $\mathcal{R}(K^\varepsilon, \mu^\varepsilon, \Delta_r, \sigma_r)$ in (5.33) can be estimated using Taylor estimates for the remainder, that is

\[
\left\| \mathcal{R}_r \right\|_{\rho} \leq C \left( \left\| \Delta_r \right\|_{\rho}^2 + |\sigma_r|^2 \right) \tag{5.42}
\]

where $C$ is a constant depending on the norms of the second derivatives of $f_{\varepsilon, \mu}$ evaluated at $K^\varepsilon$ and $\mu^\varepsilon$.

Since $f_{\varepsilon, \mu}$ is assumed to be analytic it is natural to expect the quantities

\[
\left\| D^2 f^\varepsilon \circ \mathcal{K}^\varepsilon \right\|_{\rho, \gamma_N}, \left\| D^2 f^\varepsilon \circ \mathcal{K}^\varepsilon \right\|_{\rho, \gamma_N}, \left\| D^2 f^\varepsilon \circ \mathcal{K}^\varepsilon \right\|_{\rho, \gamma_N}, \left\| D^2 f^\varepsilon \circ \mathcal{K}^\varepsilon \right\|_{\rho, \gamma_N},
\]

at the first step of the iterations. For now, we assume that $C$ is a uniform constant. In section 6, lemma 51, we give sufficient conditions on the initial error of the iteration that imply that $C$ can be taken as an uniform constant during all the iterations.

Note that (5.42) yields $\mathcal{R}_r \sim O \left( |\varepsilon|^{2N+2} \right)$. This, together with (5.37), gives (5.40). Moreover, taking sup with respect to $\varepsilon$ one obtains

\[
\left\| \mathcal{R} \right\|_{\rho-\delta, r\gamma_N} \leq C \left( \left\| \Delta_r \right\|_{\rho-\delta, r\gamma_N}^2 + \sup_{\varepsilon} \left| \sigma_r \right|^2 \right)
\]

which implies (5.43) for $\varepsilon$ small enough.
and (5.33) give the result. □

Note that the estimates above are done for the analytic functions $\Delta$ and $\sigma$. It is only left to get the respective estimates for the truncations $\Delta^{[N,2N]}$ and $\sigma^{[N,2N]}$, which are an easy consequence of the Cauchy inequalities and are given in the following propositions.

**Proposition 46.** Assume the hypothesis of lemmas 39 and 41, for any $0 < \delta < \rho$ and $0 < r < 1$ we have

$$
\mathcal{E}[K^{[\leq N]}, \mu^{[\leq N]}] + D_1 \mathcal{E}[K^{[\leq N]}, \mu^{[\leq N]}] \Delta^{[N,2N]} + D_2 \mathcal{E}[K^{[\leq N]}, \mu^{[\leq N]}] \sigma^{[N,2N]} \sim O(\|\varepsilon^{[2N+1]}\|)
$$

(5.43)

and

$$
\left\| \mathcal{E}[K^{[\leq N]}, \mu^{[\leq N]}] + D_1 \mathcal{E}[K^{[\leq N]}, \mu^{[\leq N]}] \Delta^{[N,2N]} + D_2 \mathcal{E}[K^{[\leq N]}, \mu^{[\leq N]}] \sigma^{[N,2N]} \right\|_{\rho^{-\delta}, r \gamma^N} \leq C \nu^{-3} (aN)^{2\gamma} \delta^{-(r+3\delta)} \sum_{n=1}^{N+1} \left| \frac{d}{d \delta} \mathcal{T}_\omega \right|_{\gamma^N} \| \mathcal{E}^{[N]} \|_{\rho^{-\delta}, r \gamma^N}
$$

(5.44)

**Proof.** Recalling the notation $\Delta^{[a, \infty)} = \sum_{n=a}^{\infty} \Delta_{\omega}(\theta) e^n$ we have that $\Delta^{[N,2N]} + \Delta^{[2N, \infty]} = \Delta$. Also remember that $\mathcal{E}^{[N]} = \mathcal{E}[K^{[\leq N]}, \mu^{[\leq N]}]$, then, using the linearity of the Frechet derivatives one obtains

$$
\mathcal{E}[K^{[\leq N]}, \mu^{[\leq N]}] + D_1 \mathcal{E}[K^{[\leq N]}, \mu^{[\leq N]}] \Delta^{[N,2N]} + D_2 \mathcal{E}[K^{[\leq N]}, \mu^{[\leq N]}] \sigma^{[N,2N]}
$$

$$
= \mathcal{E}[K^{[\leq N]}, \mu^{[\leq N]}] + D_1 \mathcal{E}[K^{[\leq N]}, \mu^{[\leq N]}] \Delta_{\omega}^{[2N, \infty]} - D_2 \mathcal{E}[K^{[\leq N]}, \mu^{[\leq N]}] \sigma_{\omega}^{[2N, \infty]}
$$

which implies (5.43). Moreover, using the relation above and the estimates in Proposition 44 and Proposition 42 one gets
\[
\| \mathcal{E}[K^{(\leq N)}_\varepsilon, \mu^{(\leq N)}_\varepsilon] + D_1 \mathcal{E}[K^{(\leq N)}_\varepsilon, \mu^{(\leq N)}_\varepsilon, \Delta^{(N,2N)}] + D_2 \mathcal{E}[K^{(\leq N)}_\varepsilon, \mu^{(\leq N)}_\varepsilon, \sigma^{(N,2N)}]\|_{\rho - \delta, r, N} \\
\leq \| \mathcal{E}[K^{(\leq N)}_\varepsilon, \mu^{(\leq N)}_\varepsilon] + D_1 \mathcal{E}[K^{(\leq N)}_\varepsilon, \mu^{(\leq N)}_\varepsilon, \Delta] + D_2 \mathcal{E}[K^{(\leq N)}_\varepsilon, \mu^{(\leq N)}_\varepsilon, \sigma]\|_{\rho - \delta, r, N} \\
+ \sup_{|\varepsilon| \leq \rho} \| \sigma^{(2N,\infty)}_\varepsilon \|_{\rho, \gamma, N}\]
\]

\[
\leq \frac{r^{2N+1}}{1 - r} \| E \|_{\rho, \gamma, N} + \frac{C \nu^{-d}(aN)^3 \delta^{-r+3d} r^{2N+1}}{1 - r} \| E \|_{\rho, \gamma, N}^2 \\
+ \frac{C \nu^{-1}(aN)^d \rho^{-d}(r+3d) r^{2N+1}}{1 - r} \frac{r^{2N+1}}{1 - r} \| E \|_{\rho, \gamma, N}^2 \\
\leq \frac{C \nu^{-3}(aN)^2 \delta^{-r+6d} r^{2N+1}}{1 - r} \frac{r^{2N+1}}{1 - r} \| E \|_{\rho, \gamma, N}^2
\]

**Proposition 47.** Assume the hypothesis of lemma 39 and 41, for any $0 < \delta < \rho$ and $0 < r < 1$ we have

\[
\mathcal{E} \left[ K^{(\leq N)}_\varepsilon + \Delta^{(N,2N)}_\varepsilon, \mu^{(\leq N)}_\varepsilon + \sigma^{(N,2N)}_\varepsilon \right] \sim O \left( |\varepsilon|^{2N+1} \right)
\]

and

\[
\| \mathcal{E}[K^{(\leq N)}_\varepsilon + \Delta^{(N,2N)}_\varepsilon, \mu^{(\leq N)}_\varepsilon + \sigma^{(N,2N)}_\varepsilon]\|_{\rho - \delta, r, N} \\
\leq C \nu^{-3}(aN)^2 \delta^{-r+6d} r^{2N+1} \frac{r^{2N+1}}{1 - r} \| E \|_{\rho, \gamma, N}^2
\]

where $C = C(d, \| M^{(\leq N)} \|_{\rho, \gamma, N}^{-1}, \| M^{(\leq N)} \|_{\rho, \gamma, N}^{-1}, \| \Delta^{(N,2N)} \|_{\rho, \gamma, N}^{-1}, \| D^{(N,2N)} \|_{\rho, \gamma, N}^{-1}, \mathcal{T})$, the constant $C$ also depends on the norms of the first and second derivatives of $f_{\xi, \mu}$ evaluated at $K^{(\leq N)}_\varepsilon$ and $\mu^{(\leq N)}_\varepsilon$.

**Proof.** The expansion (5.45) follows from using the same argument as in the proof of lemma 45. We also have

\[
\| \mathcal{R} \left[ K^{(\leq N)}_\varepsilon, \mu^{(\leq N)}_\varepsilon, \Delta^{(N,2N)}_\varepsilon, \sigma^{(N,2N)}_\varepsilon \right]\|_{\rho - \delta, r, N} \\
\leq C \left( \| \Delta^{(N,2N)}_\varepsilon \|_{\rho - \delta, r, N}^2 + \sup_{|\varepsilon| \leq \rho} \| \sigma^{(N,2N)}_\varepsilon \|_{\rho, \gamma, N}^2 \right)
\]
\[ \leq C \left( \nu^{-\theta}(aN)^{3\theta} \delta^{-(2r+6d)} \frac{r^{2N+2}}{(1-r^2)^2} \left\| E^N \right\|_{\rho_0, \delta_N}^2 + \nu^{-2}(aN)^{2\theta} \rho^{-2d} \frac{r^{2N+2}}{(1-r^2)^2} \left\| E^N \right\|_{\rho_0, \delta_N}^2 \right) . \]

Combining this estimate with (5.44) in proposition 46 one gets (5.46).

\[ \square \]

6. Iteration of the quasi-Newton method

We start this section giving the choice of parameters which quantify the loss of regularity at any step of the quasi Newton method. Lemma 51 will guarantee that the Newton method is well defined at any step. We note that we have loss of domain in both the variable on the torus, \( \theta \), and the variable of the perturbation, \( \varepsilon \). In contrast with the regular KAM theory we end up losing much more domain in \( \varepsilon \), so that at the end we do not have any \( \varepsilon \) domain.

6.1. The iterative procedure

We denote by \( h \in \mathbb{N} \) the number of steps of the quasi Newton method. We consider

\[ \delta_h := \frac{\rho_0}{2^h+2} \text{ and } \rho_{h+1} := \rho_h - \delta_h \geq \rho_0 \frac{2}{2^h} \text{ for } h \geq 1, \quad (6.1) \]

where \( \rho_h \) denotes the radius of analyticity in the variable \( \theta \) at step \( h \), that is, at step \( h \) we will be considering functions in the space \( A_{\rho_h} \). Note that \( \rho_0 = \rho \) can be the one given in theorem 19.

Since at any step we double the number of coefficients of the Lindstedt expansions, we have,

\[ N_h := 2^h N_0 \quad (6.2) \]

and

\[ \gamma_h := \gamma_N = \left( \frac{\nu}{2} \right)^{1/\alpha} \frac{1}{(aN)^{\gamma/\alpha}} = \left( \frac{\nu}{2} \right)^{1/\alpha} \frac{1}{(a2^h N_0)^{\gamma/\alpha}}, \quad (6.3) \]

where \( \alpha \in \mathbb{N} \) is the exponent in \( \lambda(\varepsilon) = 1 - \varepsilon^\alpha \), \( a \in \mathbb{N} \), and \( N_0 \in \mathbb{N} \) is a fixed constant to be chosen later. Note that \( \gamma_h \) is the radius of the domain of analyticity in the variable \( \varepsilon \) at step \( h \), that is, at step \( h \) we will be considering functions in the space \( A_{\gamma_h} \). Also note that

\[ \gamma_{h+1} = 2^{-\tau/\gamma} \gamma_h. \quad (6.4) \]

Denoting \( K_0 := K^{[\leq N_0]} \) and \( \mu_0 := \mu^{[\leq N_0]} \), for \( h \geq 1 \) we have

\[ K_h := K^{[\leq N_h]} + \Delta^{[N_h, N_{h+1}]} + \ldots + \Delta^{[N_{h+1}, N_h]} \]

\[ \mu_h := \mu^{[\leq N_h]} + \sigma^{[N_h, N_{h+1}]} + \ldots + \sigma^{[N_{h+1}, N_h]}, \quad (6.5) \]

Furthermore, denoting

\[ \Delta_h := \Delta^{[N_h, N_{h+1}]} \text{ and } \sigma_h := \sigma^{[N_h, N_{h+1}]} \text{ for } h \geq 0 \]

we have that, for \( h \geq 0 \)

\[ K_{h+1} = K_h + \Delta_h \text{ and } \mu_{h+1} = \mu_h + \sigma_h, \quad (6.7) \]
Finally, denote also
\[ e_h := \|E(K_h, \mu_h)\|_{\rho_h, \tilde{\gamma}_h} = \|E_N\|_{\rho_h, \tilde{\gamma}_h} \] (6.8)
\[ d_h := \|\Delta_h\|_{\rho_{h+1, 5h+1}} \] (6.9)
\[ v_h := \|D\Delta_h\|_{\rho_{h+1, 5h+1}} \] (6.10)
\[ s_h := \sup_{|t| \leq \tilde{\gamma}_h+1} |\sigma_h(t)|. \] (6.11)

**Remark 48.** We emphasize the dependence of \( \tilde{\gamma}_h \) in \( N_h \), note that \( \tilde{\gamma}_h \to 0 \) as \( N_h \to \infty \) \((h \to \infty)\). This implies that this quasi Newton method will not converge in any Banach space \( \mathcal{A}_{\rho_0, \tilde{\gamma}_0} \), because the domains in \( \varepsilon \) shrink to 0, however, at each step we get estimates in balls with positive radius, \( \tilde{\gamma}_h \). An analysis of these bounds will provide us with estimates of the coefficients of the expansion. Note also that to start with \( e_0 \ll 1 \) we require \( N_0 \) sufficiently large in the formal power series in theorem 19.

Note that with this new notation the estimates in corollary 42 can be written as
\[ d_h \leq \bar{C}_h \nu^{-3} (aN_h)^{2\tau} \delta_h^{-\tau+3d} \left( \frac{1}{2\gamma} \right)^{N_h} e_h \] (6.12)
\[ v_h \leq \bar{C}_h \nu^{-3} (aN_h)^{2\tau} \delta_h^{-\tau+3d+1} \left( \frac{1}{2\gamma} \right)^{N_h} e_h \] (6.13)
\[ s_h \leq \bar{C}_h \nu^{-1} (aN_h)^{\gamma} \delta_h^{-d} \left( \frac{1}{2\gamma} \right)^{N_h} e_h, \] (6.14)

where \( \bar{C}_h \) is an explicit constant depending in a polynomial manner on \( \|M_h\|_{\rho_0, \tilde{\gamma}_0}, \|M_h^{-1}\|_{\rho_0, \tilde{\gamma}_0}, \|N_h\|_{\rho_0, \tilde{\gamma}_0}, \|D\mathcal{K}_h\|_{\rho_0, \tilde{\gamma}_0}, \) and \( T_h \). Moreover, the nonlinear estimate (5.46) given in proposition 47 implies
\[ e_{h+1} \leq \bar{C}_h (aN_h)^{\tau} \delta_h^{-(2\tau+6d)} \left( \frac{1}{2\gamma} \right)^{N_h} \left( e_h + e_h^2 \right), \] (6.15)

where \( \bar{C}_h \) is a constant which also depends explicitly on \( \|M_h\|_{\rho_0, \tilde{\gamma}_0}, \|M_h^{-1}\|_{\rho_0, \tilde{\gamma}_0}, \|N_h\|_{\rho_0, \tilde{\gamma}_0}, \|D\mathcal{K}_h\|_{\rho_0, \tilde{\gamma}_0}, \) and \( T_h \).

**Remark 49.** In the following we will denote \( C \) a constant depending on \( \nu, \tau, d, \xi, \rho_0, |J^{-1}| \), and that is a polynomial in \( \|M_h\|_{\rho_0, \tilde{\gamma}_0}, \|M_h^{-1}\|_{\rho_0, \tilde{\gamma}_0}, \|N_h\|_{\rho_0, \tilde{\gamma}_0}, \|D\mathcal{K}_h\|_{\rho_0, \tilde{\gamma}_0}, \) and \( T_h \). We will also denote
\[ C_h = \max \left( \bar{C}_h, \tilde{C}_h \right). \]

In lemma 51, we give smallness conditions so that \( C_h \leq C \) for every \( h \geq 0 \). Since we are working with expansions near to \((K^{\gamma(\leq \alpha)}), \mu^{\gamma(\leq \alpha)}\) it is natural to expect that the quantities \( \|M_h\|_{\rho_0, \tilde{\gamma}_0}, \|M_h^{-1}\|_{\rho_0, \tilde{\gamma}_0}, \|N_h\|_{\rho_0, \tilde{\gamma}_0}, \|D\mathcal{K}_h\|_{\rho_0, \tilde{\gamma}_0}, \) and \( T_h \) will be close to \( \|M_h\|_{\rho_0, \tilde{\gamma}_0}, \|M_h^{-1}\|_{\rho_0, \tilde{\gamma}_0}, \|N_h\|_{\rho_0, \tilde{\gamma}_0}, \|D\mathcal{K}_h\|_{\rho_0, \tilde{\gamma}_0}, \) and \( T_h \), respectively. For now, we assume that \( C \) is large enough, for instance \( C > 2C_0 \). Here \( M_h = M^{\gamma(\leq \alpha)}, N_h = N^{\gamma(\leq \alpha)}, \) and \( T_h = T^{\alpha} \) as in (4.8), (4.10) and (5.20).
Considering this uniform constant $C$ on (6.15), and taking $N_0$ sufficiently large, yields $e_h < 1$ for any $h > 0$, and inequality (6.15) implies
\[ e_{h+1} \leq C N^{-0}(aN_0)^{4\tau} \sigma_0^{-((2+6d)2^{-4(4+12d)}(2^h)^{6(2^h)+6d} \left( \frac{1}{2^{7/\alpha}} \right)^{2\delta N_0}} e_h. \]  

(6.16)

**Remark 50.** Due to remark 49 and the definitions of $\delta_h, \rho_h, N_h,$ and $\gamma_h$, the inequality (6.16) can be rewritten as
\[ e_{h+1} \leq C N^{-0}(aN_0)^{4\tau} \rho_0^{-((2+6d)2^{-4(4+12d)}(2^h)^{6(2^h)+6d} \left( \frac{1}{2^{7/\alpha}} \right)^{2\delta N_0}} e_h \]
or
\[ e_{h+1} \leq CDB^h \rho^{2\delta N_0} e_h, \]  

(6.17)

where
\[ D = \nu^{-6}(aN_0)^{4\tau} \rho_0^{-((2+6d)2^{-4(4+12d)}(2^h)^{6(2^h)+6d} \left( \frac{1}{2^{7/\alpha}} \right)^{2\delta N_0}} r^{-2^{-7/\alpha}} \text{ and } B = 2^{6(2^h)+6d}. \]

**Lemma 51.** Assume that $2^{3(\tau+3d)+1} CDB^0 N_0 \leq \frac{1}{2}, B^0 N_0 < 1, N_0^2 e_0 \ll 1$, and $C N^{-3}(aN_0)^{3\tau} \rho_0^{-((\tau+3d)+1)2^{(2r+6d+2)} e_0} \ll 1$. Then, for all integers $h \geq 0$ the following properties hold:

- **[(p1; $h$)]**
  \[ \|K_h - K_0\|_{p_0, \gamma_h} \leq \ell_k N_0^2 e_0 < \xi \]
  \[ \sup |\mu_h - \mu_0| \leq \ell_\mu N_0^2 e_0 < \xi \]
  with $\ell_k \equiv C N^{-3}(aN_0)^{3\tau} \rho_0^{-((\tau+3d)+1)2^{(2r+6d+2)} e_0}$ and $\ell_\mu \equiv C N^{-1} a^2 \rho_0^{-d}$

- **[(p2; $h$)]**
  \[ e_h \leq (CD)^2 B^2 \ell^{(2^h+1)N_0} e_0 \]

- **[(p3; $h$)]**
  \[ C_h \leq C \]

**Remark 52.** Note that by (3.12) we have $e_0 \sim O(N_0^{-\tau/\alpha})$, due to the fact that we estimate $e_0$ in a ball with radius $\tilde{\gamma}_0 \sim O(N_0^{-\tau/\alpha})$. So the assumptions on the smallness of $N_0 e_0$ are satisfied.

**Proof.** Note that (p1; 0), (p2; 0), and (p3; 0) are trivial.

Let us now prove (p1; $H + 1$), (p2; $H + 1$), and (p3; $H + 1$) assuming they are true for $h = 1, 2, \ldots, H$. Noticing that $2^j \leq 2^{j+1} - 1$, for any $j \geq 0$, and assuming that $N_0$ is large enough such that $2^{3(\tau+3d)+1} CDB^0 N_0 \leq \frac{1}{2}$ and $B^0 N_0 < 1$, we have
\[ \| K_{H+1} - K_0 \|_{pH+1,1} = \| \Delta^{(N_0, N_0)} \cdots + \Delta^{(N_H, N_H+1)} \|_{pH+1,1} \leq \sum_{j=0}^{H} d_j \leq \sum_{j=0}^{H} C_j \nu^{-3} (a N_j)^{2^r} \delta_j^{-(\tau+3d)} \rho_j e_j \]

\[ \leq C \nu^{-3} (a N_0)^{2^r} \rho_0^{-(\tau+3d)} 2^{r+6d} \]

\[ \leq C \nu^{-3} (a N_0)^{2^r} \rho_0^{-(\tau+3d)} 2^{r+6d} \]

\[ \times \sum_{j=0}^{H} 2^{3(d+7j)} j^{2j/2} (\sigma^{(N_H, N_H+1)N_0}) e_j \]

\[ \leq C \nu^{-3} (a N_0)^{2^r} \rho_0^{-(\tau+3d)} 2^{r+6d} \]

\[ \times \sum_{j=0}^{H} 2^{3(d+7j)} (CD) B_j^{r^2} \rho^{(r+1)N_0} e_0 \]

\[ \leq C \nu^{-3} (a N_0)^{2^r} \rho_0^{-(\tau+3d)} 2^{r+6d} \]

\[ \times \sum_{j=0}^{H} (2^{3(d+7j)} CDBr^{N_0})^j \]

\[ \leq C \nu^{-3} (a N_0)^{2^r} \rho_0^{-(\tau+3d)} 2^{r+6d} \]

\[ \leq C \nu^{-3} (a N_0)^{2^r} \rho_0^{-(\tau+3d)} 2^{r+6d} e_0 \]

\[ \leq \ell K N_0^{2^r} e_0. \]

Similarly,

\[ \sup_{|\epsilon| \leq \tilde{\gamma}_{H+1}} |\mu_{H+1} - \mu_0| = \sup_{|\epsilon| \leq \tilde{\gamma}_{H+1}} |\sigma^{(N_0, N_0)} + \cdots + \sigma^{(N_H, N_H+1)}| \]

\[ \leq \sum_{j=0}^{H} s_j \leq \sum_{j=0}^{H} C_j \nu^{-1} (a N_j)^{2^r} \delta_j^{-(\tau+3d)} \rho_j e_j \]

\[ \leq \sum_{j=0}^{H} C \nu^{-1} (a N_0)^{2^r} \rho_0^{-(\tau+3d)} 2^{2j+2d} \rho^{N_0} \]

\[ \times (CD) B_j^{r^2} \rho^{(r+1)N_0} e_0 \]

\[ \leq C \nu^{-1} (a N_0)^{2^r} \rho_0^{-(\tau+3d)} 2^{2d} \sum_{j=0}^{H} (2^{r+d})^j \]

\[ \times (CD) B_j^{r^2} \rho^{(r+1)N_0} e_0 \]
Thus, taking $N_0$ large enough, which makes $\epsilon_0$ small, we get $\ell_p N_0^2 \epsilon_0 < \xi$ and $\ell_p N_0^2 \epsilon_0 < \xi$.

Since $(p_3; H + 1)$ is true, we use the estimate (6.17) given in remark 50, which is a consequence of the nonlinear estimates given in proposition 47, that is

$$e_{h+1} = \| \hat{\xi}(K_h + \Delta_h, \mu_h + \sigma_h) \|_{\rho_{h+1}, \gamma_{h+1}} \leq CDB^h r^{2h} N_0 \epsilon_h, \quad (6.18)$$

where $D$, $B$, and $r$ are as in remark 50. This yields,

$$e_{h+1} \leq CDB^h r^{2h} N_0 \epsilon_h$$

which yields $(p_3; H + 1)$.

In order to prove $(p_3; H + 1)$ note that

$$\| N_0 - N_0 \|_{\rho_{h+1}, \gamma_{h+1}} \leq \bar{C} \| DK_h - DK_0 \|_{\rho_{h+1}, \gamma_{h+1}} \quad (6.19)$$

$$\| M_h - M_0 \|_{\rho_{h+1}, \gamma_{h+1}} \leq \bar{C} \| DK_h - DK_0 \|_{\rho_{h+1}, \gamma_{h+1}} \quad (6.20)$$

$$\| M_h^{-1} - M_0^{-1} \|_{\rho_{h+1}, \gamma_{h+1}} \leq \bar{C} \| DK_h - DK_0 \|_{\rho_{h+1}, \gamma_{h+1}} \quad (6.21)$$

$$\| T_h - T_0 \|_{\rho_{h+1}, \gamma_{h+1}} \leq \bar{C} \| DK_h - DK_0 \|_{\rho_{h+1}, \gamma_{h+1}} \quad (6.22)$$

where $\bar{C}$ is a uniform constant. The above inequalities come from the fact that $M_h, N_h, \text{ and } T_h$ are algebraic expressions of $DK_h, Df_{\rho_0}$, and $\mu_{\rho_0}$; see (4.8), (4.10), (4.9), (5.20). Then,

$$\| DK_{h+1} - DK_0 \|_{\rho_{h+1}, \gamma_{h+1}} = \| D\Delta^{(N_0 N_1)} + \cdots + D\Delta^{(N_0 N_{H+1})} \|_{H+1, \gamma_{H+1}}$$

$$\leq \sum_{j=0}^{H} d_j \leq \sum_{j=0}^{H} C \nu^{-3} (a N_0)^{2r} \rho_0^{-(r+3d+1)} \rho_{h+1} \epsilon_j$$

$$\leq \sum_{j=0}^{H} C \nu^{-3} (a^2 N_0)^{2r} \rho_0^{-(r+3d+1)}$$

$$\times 2^{2r+6d+2}(r+3d+1)j^2 N_0 \epsilon_j$$

$$2464$$
given 

\[ K \]

for any 

Note that theorem 19 assures the existence of the Lindstedt series satisfying (6.2).

Proof. 

For this proof we inherit all the notation introduced throughout this section.

6.2. Proof of theorem B

For this proof we inherit all the notation introduced throughout this section.

Proof. 

Note that theorem 19 assures the existence of the Lindstedt series satisfying (6.2).

That is, we can construct the unique formal powerseries as follows

\[
\leq C\nu^{-3}(aN_0)^{2(\nu-3d+3)}\rho_0^{-2r+6d+2} \times \sum_{j=0}^{H} 2^{3d+3j+1}J_{j/2N_0} \times \left((CD)^{j/2}r_j^{2^{j-1}N_0}\epsilon_0\right)
\]

\[
\leq C\nu^{-3}(aN_0)^{2(\nu-3d+3)}\rho_0^{-2r+6d+2} \times \sum_{j=0}^{H} 2^{3d+3j+1}(CD)^{j/2}r_j^{2^{j-1}N_0}\epsilon_0
\]

\[
\leq C\nu^{-3}(aN_0)^{2(\nu-3d+3)}\rho_0^{-2r+6d+2} \times \sum_{j=0}^{H} (2^{3d+3j+1}CDBr_0)^{j/2}
\]

\[
\leq C\nu^{-3}(aN_0)^{2(\nu-3d+3)}\rho_0^{-2r+6d+2}\epsilon_0,
\]

where the sum is bounded as in the previous estimates. Taking \( \epsilon_0 \) small enough, such that 

\[ C\nu^{-3}(aN_0)^{2(\nu-3d+3)}\rho_0^{-2r+6d+2}\epsilon_0 \ll 1, \]

we are able to verify (p3; H + 1) because \( C_{H+1} \) is an algebraic expression of \( M_H, N_H, \) and \( T_H; \) and taking \( C \geq 2C_0, \) for example.

\[ \square \]

6.2. Proof of theorem B

For this proof we inherit all the notation introduced throughout this section.

\[ \square \]

Proof. 

Note that theorem 19 assures the existence of the Lindstedt series satisfying (6.2).

That is, given \( K_0 \in A, \) and \( \mu_0 \in A \subseteq C \) satisfying \( f_{0,\mu_0} \circ K_0 = K_0 \circ T_N, \) and \( H_N, \) there exists 

\[ \rho_0 < \rho \]

and power expansions \( K^{[\leq N]} \) and \( \mu^{[\leq N]} \) such that

\[
\left\| f_{0,\mu_0}^{[\leq N]} \circ K^{[\leq N]} \circ T_N \right\|_{\rho_0} \leq C_N|\nu|^{N+1}
\]

for any \( N \geq 0. \) This expansion is unique under the normalization condition (3.3).

If \( K^{[\leq N]} \) and \( \mu^{[\leq N]} \) satisfy hypothesis HTP1 and HTP2 then, we can choose \( N_0 \) such that 

\[ K^{[\leq N_0]} \]

and \( \mu^{[\leq N_0]} \) satisfy the hypothesis of lemmas 39 and 41. Also, \( N_0 \) needs to be large enough such that 

\[ 2^{3r+3d+1}CDBr_0 \leq \frac{1}{2}, B^{r_0} < 1, \epsilon_0 N_0^2 \epsilon_0 < \xi, \epsilon_0 N_0 \epsilon_0 \ll 1, \]

\[ \text{then lemma 51 can be applied and this allows us to iterate the quasi Newton method described} \]

in algorithm 35. That is, we can construct the unique formal power series as follows

\[ \square \]
\[
K_{\varepsilon}^{(\leq N_0)} + \Delta_k^{(N_0,2N_0)} + \Delta_k^{(2N_0,2^2N_0)} + \ldots + \Delta_k^{(2^k N_0,2^{k+1} N_0)} + \ldots \tag{6.23}
\]

\[
\mu_{\varepsilon}^{(\leq N_0)} + \mu_{\varepsilon}^{(N_0,2N_0)} + \mu_{\varepsilon}^{(2N_0,2^2N_0)} + \ldots + \mu_{\varepsilon}^{(2^k N_0,2^{k+1} N_0)} + \ldots \tag{6.24}
\]

Note that by definition of \(\tilde{\gamma}_k\) we will have \(\tilde{\gamma}_k = r^k \tilde{\gamma}_0\), where \(r = 2^{-\tau/\alpha}\) and \(\tilde{\gamma}_0 = 2^{-1/\alpha} \nu^{1/\alpha} (aN_0)^{-\tau/\alpha}\), see (6.4). Before giving the detailed computations, note that \(\tilde{\gamma}_k = (2^{-1} \nu a^{-\tau})^{1/\alpha} (2^k N_0)^{-\tau/\alpha}\) and if \(n \in (2^k N_0, 2^{k+1} N_0]\cap \mathbb{N}\) then, using the notation \(c = (2^{-1} \nu a^{-\tau})^{1/\alpha}\), one has

\[
(\tilde{\gamma}_k)^{-n} = c^n (2^k N_0)^{(\tau/\alpha)n} \approx c^n n^{(\tau/\alpha)n}. \tag{6.25}
\]

Recalling the notation \(\Delta_h := \Delta_k^{(2^k N_0, 2^{k+1} N_0)}\), we also note that (6.12), and lemma 51 imply that

\[
\|\Delta_h\|_{h_0+1, \tilde{\gamma}_h+1} \leq \hat{C}_h \nu^{-3} (aN_0)^2 \tau \delta_h \tau (r+3d) \rho_0 \varepsilon_h
\]

\[
\leq C \nu^{-3} (2^k N_0)^{2r} \left(\rho_0^{-1} 2^{k+2}\right)^{\tau (r+3d)} \rho_0 \varepsilon_h
\]

\[
\times (CD)^h B^2 \tau (2^{k+1} N_0) \varepsilon_0
\]

that is

\[
\|\Delta_h\|_{h_0+1, \tilde{\gamma}_h+1} \leq L c_0^2 r^{2^{k+1} N_0}, \tag{6.26}
\]

where \(c_0 = 2^{3r+3d} C D B > 1\) and the constant term in (6.26) given by \(L = C \nu^{-3} (aN_0)^{2r} (\rho_0^{-1} 2^{k+2})^{\tau (r+3d)} \rho_0 \varepsilon_0\).

Using the observations above together with Cauchy estimates yield the Gevrey estimates. More precisely, note that for \(n \in (2^k N_0, 2^{k+1} N_0]\cap \mathbb{N}\), the coefficient \(K_n\) of the expansion \(\Delta_h\) belongs to the correction \(\Delta_{h}\), then using Cauchy estimates and (6.26), we obtain

\[
\|K_n\|_{2n} \leq (\tilde{\gamma}_h+1)\quad \|\Delta_h\|_{2n, \tilde{\gamma}_h+1}
\]

\[
\leq (\tilde{\gamma}_h+1)^{-n} \|\Delta_h\|_{h_0+1, \tilde{\gamma}_h+1}
\]

\[
= \left(r^{h+1} \tau \rho_0^{-1} \right)^{-n} L c_0^2 \rho_0 \varepsilon_0 \tau^{2^{k+1} N_0}
\]

\[
= \left(r^{h+1} \tau \rho_0^{-1} \right)^{-n} L c_0^2 \rho_0 \varepsilon_0 \tau^{2^{k+1} N_0}
\]

\[
\leq \left(c_0^2 \rho_0 \varepsilon_0 \tau^{2^{k+1} N_0}\right)^{2n \varepsilon_0 \tau^{2^{k+1} N_0}}
\]

\[
\leq LF^n (2^k N_0)^{\tau n / \alpha}
\]

\[
\leq LF^n n^{(\tau / \alpha)n},
\]

where \(F = c_0^2 \rho_0 \varepsilon_0 \tau^{2^{k+1} N_0}\). The estimates for \(\mu_n\) are obtained in a similar way. \(\square\)
6.3. Proof of theorem C

Proof. Denoting $K_{\varepsilon}^{[n]} := \sum_{j=0}^{n} K_{\varepsilon}^{j} / n$ and $E_{\varepsilon}^{n} := f_{\varepsilon,\mu_{\varepsilon}} \circ K_{\varepsilon}^{[n]} - K_{\varepsilon}^{[n]} \circ T_{\omega}$.

If for all $\varepsilon \in \mathcal{G}$ such that $|\varepsilon| \leq \gamma_{0}$, $\|E_{\varepsilon}^{n}\|_{\rho_{0}} \leq C (\nu^{-1} \nu(\varepsilon, \omega, \tau))^{2}\delta^{4(r+d)}$ (6.27)

then, by theorem 14 in [CCdlL17] one obtains that there exist unique functions $K_{\varepsilon}$ and $\mu_{\varepsilon}$, defined in $\varepsilon \in \mathcal{G} \cap B_{\gamma_{h}}(0)$ satisfying,

$\|K_{\varepsilon}^{[n]} - K_{\varepsilon}\|_{\rho_{0}} \leq \tilde{C} \nu(\varepsilon, \omega, \tau)^{-1}\delta^{-2(r+d)}\|E_{\varepsilon}^{n}\|_{\rho_{0}}$ (6.28)

and

$|\mu_{\varepsilon}^{[n]} - \mu_{\varepsilon}| \leq \tilde{C}\|E_{\varepsilon}^{n}\|_{\rho_{0}}\gamma_{h+1}$ (6.29)

where (6.29) follows from the fact that $\nu(\varepsilon, \omega, \tau)$ is lower semi-continuous and that $\nu(0; \omega, \tau) = \nu^{-1}$, see remark 22. The Cauchy estimates in corollary 6 imply that, if $n \in (2^{h}N_{0}, 2^{h+1}N_{0}] \cap \mathbb{N}$ then

$\|E_{\varepsilon}^{n}\|_{\rho_{0}} \leq C r^{h} \|E_{\rho_{0}}^{2^{h}N_{0}}\|_{\rho_{0}}$ (6.31)

then, using the notation $e_{h} := \|E_{\rho_{0}}^{2^{h}N_{0}}\|_{\rho_{0}} \gamma_{h}$ and the estimate $(p, h)$ in lemma 51, that is $e_{h} \leq (CD)^{h}B^{h}r^{2^{h+1}N_{0}}\|E_{\rho_{0}}^{2^{h}N_{0}}\|_{\rho_{0}}$, one has from (6.29) that for any $\varepsilon \in \mathcal{G} \cap B_{\gamma_{h+1}}(0)$

$\|K_{\varepsilon}^{[n]} - K_{\varepsilon}\|_{\rho_{0}} \leq \tilde{C} \nu^{-1}\delta^{-2(r+d)}\|E_{\varepsilon}^{n}\|_{\rho_{0}}\gamma_{h}$

and

$|\mu_{\varepsilon}^{[n]} - \mu_{\varepsilon}| \leq \tilde{C}\|E_{\varepsilon}^{n}\|_{\rho_{0}}\gamma_{h}$

$\leq \tilde{C}(CD)^{h}B^{h}r^{2^{h}N_{0}}.$

Acknowledgments

The authors would like to thank T M-Seara, I Baldomá, and V Naudot for many suggestions and discussions. We also like to thank to the anonymous reviewers for their valuable comments and suggestions.

Appendix A. The case of the dissipative standard map of theorem A

A.1. Verifying HTP1 and HTP2 for the dissipative standard map

Consider the dissipative standard map $f_{\varepsilon,\mu_{\varepsilon}} : T \times \mathbb{R} \to T \times \mathbb{R}$ given by
where \( V(x) \) is a trigonometric polynomial. In this section we verify that maps like (A.1) satisfy HTP1 and HTP2 of theorem B. For the sake of simplicity in the exposition we do it for the case \( \lambda(\varepsilon) = 1 - \varepsilon^3 \). The general case for \( \alpha \in \mathbb{N} \) is done by very similar computations, fixing the value of \( \alpha = 3 \) allows an easy analysis of the Lindstedt series.

Note that one has \( f_{\varepsilon, \mu} \Omega = \lambda(\varepsilon)\Omega \) for the symplectic form \( \Omega_{(x,y)} = dx \wedge dy \), so it is conformally symplectic. One can write the map as

\[
x_{n+1} = x_n + y_{n+1}
\]
\[
y_{n+1} = \lambda(\varepsilon)y_n + \mu_\varepsilon - \varepsilon V(x_n)
\]
equivalently

\[
x_{n+1} - (1 + \lambda(\varepsilon))x_n + \lambda(\varepsilon)x_{n-1} - \mu_\varepsilon + \varepsilon V(x_n) = 0. \tag{A.2}
\]

Considering a parametric representation of the variable \( x_n \in \mathbb{T} \) as \( x_n = \theta_n + u_\varepsilon(\theta_n), \theta_n \in \mathbb{T} \);

where \( u_\varepsilon : \mathbb{T} \rightarrow \mathbb{R} \) is a one-periodic function and assuming that \( \theta_n \) varies linearly, i.e. \( \theta_{n+1} = \theta_n + \omega \), then, (A.2) becomes

\[
u_\varepsilon(\theta + \omega) - (1 + \lambda(\varepsilon))u_\varepsilon(\theta) + \lambda(\varepsilon)u_\varepsilon(\theta - \omega) + (1 - \lambda(\varepsilon))\omega - \mu_\varepsilon + \varepsilon V(\theta + u_\varepsilon(\theta)) = 0. \tag{A.3}
\]

If \( u_\varepsilon \) satisfies (A.3) it is easy to check that \( K_\varepsilon : \mathbb{T} \rightarrow \mathbb{T} \times \mathbb{R} \), given by

\[
K_\varepsilon(\theta) = \left( \frac{\theta + u_\varepsilon(\theta)}{\omega + u_\varepsilon(\theta) - u_\varepsilon(\theta - \omega)} \right),
\]
satisfies \( f_{\varepsilon, \mu} \circ K_\varepsilon(\theta) = K(\theta + \omega) \). Therefore, the problem of finding Lindstedt series for quasiperiodic orbits for the map \( f_{\varepsilon, \mu} \) is equivalent to find asympototic power series to a solution, \( (u_\varepsilon, \mu_\varepsilon) \), of (A.3).

Using \( \lambda(\varepsilon) = 1 - \varepsilon^3 \), equation (A.3) becomes

\[
u_\varepsilon(\theta + \omega) - (2 - \varepsilon^3)u_\varepsilon(\theta) + (1 - \varepsilon^3)u_\varepsilon(\theta - \omega) + \varepsilon^3 \omega - \mu_\varepsilon + \varepsilon V(\theta + u_\varepsilon(\theta)) = 0. \tag{A.4}
\]

Introducing the operator

\[
L_\varepsilon u(\theta) = u(\theta + \omega) - 2u(\theta) + u(\theta - \omega),
\]

and expanding in power series on \( \varepsilon \), i.e. \( u_\varepsilon(\theta) = \sum_{n=0}^{\infty} u_n(\theta)\varepsilon^n \) and \( \mu_\varepsilon = \sum_{n=0}^{\infty} \mu_n \varepsilon^n \), equation (A.4) becomes

\[
\sum_{k=0}^{2} (L_\varepsilon u_k(\theta) - \mu_k) \varepsilon^k - (L_\varepsilon u_3(\theta) - \mu_3 + u_0(\theta) - u_0(\theta - \omega) - \omega) \varepsilon^3
\]
\[
+ \sum_{k=4}^{\infty} (L_\varepsilon u_k(\theta) - \mu_k + u_{k-3}(\theta) - u_{k-3}(\theta - \omega)) \varepsilon^k = -\sum_{k=1}^{\infty} S_{k-1}(\theta) \varepsilon^k. \tag{A.5}
\]

**Remark 53.** When \( V(\theta) \) is a trigonometric polynomial, the coefficients \( S_n \) can be computed as follows. Note that \( V_k(\theta) = f_{\varepsilon} e^{2\pi i k \theta} \) satisfies the relation
\[
\frac{d}{d\varepsilon}V_k(\theta + u_\varepsilon(\theta)) = 2\pi i k \frac{d}{d\varepsilon}u_\varepsilon(\theta)V_k(\theta + u_\varepsilon(\theta)).
\] (A.6)

Thus, considering
\[
V_k(\theta + u_\varepsilon(\theta)) = \sum_{n=0}^{\infty} S_n^k(\theta)\varepsilon^n
\]
and (A.6) the coefficients \(S_n^k\) satisfy the following relation
\[
(n+1)S_{n+1}^k = \sum_{\ell=0}^{n} 2\pi i k (\ell+1)u_{\varepsilon+1}\varepsilon^\ell S_n^{k-\ell},
\] (A.7)
and \(S_0^k(\theta) = \hat{V}_k e^{2\pi i k \theta}\). Furthermore, if \(V(\theta) = \sum_{|k| \leq a} \hat{V}_k e^{2\pi i k \theta} = \sum_{|k| \leq a} V_k(\theta)\) is a trigonometric polynomial of degree \(a\), considering
\[
V(\theta + u_\varepsilon(\theta)) = \sum_{n=0}^{\infty} S_n(\theta)\varepsilon^n,
\]
the coefficients \(S_n(\theta)\) are given by
\[
S_n(\theta) = \sum_{|k| \leq a} S_n^k(\theta),
\]
where \(S_n^k\) is given by (A.7).

**Remark 54.** Note that if \(\eta\) is a trigonometric polynomial and \(\varphi\) is a solution of the equation \(L_\omega \varphi = \eta\) then, \(\varphi\) is a trigonometric polynomial of the same degree as \(\eta\). This is due to the fact that the Fourier coefficients of \(\varphi\) satisfy \(\hat{\varphi}_k = \frac{1}{2\pi i k} \hat{\eta}_k\). Note that the equation \(L_\omega \varphi = \eta\) has a solution if \(\int \varphi(\theta) d\theta = 0\), and this solution is unique if we impose the normalization \(\int \varphi(\theta) d\theta = 0\).

**Proposition 55.** If \(V(\theta)\), in (A.1), is a trigonometric polynomial of degree \(a\), then \(u_\varepsilon(\theta)\) is a trigonometric polynomial of degree \(an\). Furthermore, \(S_{n-1}(\theta)\) is a trigonometric polynomial of degree \(an\).

**Proof.** Equating the terms of same order in equation (A.5) one gets that for order zero \(\mu_0 = 0\) and \(u_\varepsilon(\theta) \equiv 0\). For order 1 we have,
\[
L_\omega u_1(\theta) - \mu_1 = -S_0(\theta).
\]
So, taking \(\mu_1 = 0\), \(u_1\) becomes a trigonometric polynomial of degree \(a\), because \(S_0(\theta) = V(\theta)\). Now, for order 2 we have
\[
L_\omega u_2(\theta) - \mu_2 = -S_1(\theta),
\]
if \(\mu_2 = 0\) the right-hand side is \(S_1(\theta) = \sum_{|k| \leq a} S_1^k(\theta) = 2\pi i u_1(\theta) \sum_{|k| \leq a} k S_0^k(\theta)\) which is a trigonometric polynomial of degree \(2a\), thus \(u_2\) is a trig polynomial of degree \(2a\). For order three we have
\[
L_\omega u_3(\theta) - \mu_3 + \omega = -S_2(\theta),
\]
here we take $\mu_3 = \omega$ and $\nu_1$ is a trig polynomial of degree $3a$ because

$$S_3(\theta) = \sum_{|k| \leq a} S^k_3(\theta) = \pi i u_1(\theta) \sum_{|k| \leq a} k S^k_1(\theta) + 2 \pi i u_2(\theta) \sum_{|k| \leq a} k S^k_0(\theta)$$

is of degree $3a$; then $u_3(\theta)$ is of degree $3a$. Finally, for $n > 4$, assume the claim is valid for any $m < n$ then, the equation of order $n$ is

$$L_m u_n(\theta) = \mu_m - u_{n-3}(\theta) + u_{n-3}(\theta - \omega) - S_{n-1}(\theta).$$

So, taking $\mu_m = \int S_{n-1}(\theta) d\theta$, $u_n$ can be found and has degree $an$ since, $S_{n-1} = \sum_{|k| \leq n} S^k_{n-1}$ and each $S^k_{n-1}$ has degree $an$ due to (A.7). Note $u_{n-3}$ has degree $(n - 3)a$.

**Corollary 56.** If $V(\theta)$, in (A.1), is a trigonometric polynomial of degree $a$, then for any fixed $\varepsilon$ the sum $\sum_{n=0}^{N} u_n(\varepsilon \cdot n)$ is a trigonometric polynomial of degree $aN$ in $\theta$.

Note that in this case

$$K_\varepsilon^{[\leq N]}(\theta) = \begin{pmatrix} \theta + \sum_{n=0}^{N} u_n(\theta) \varepsilon^n \\ \omega + \sum_{n=0}^{N} (u_n(\theta) - u_n(\theta - \omega)) \varepsilon^n \end{pmatrix},$$

and using equation (A.5) we have

$$E_\varepsilon^N(\theta) := f_{\varepsilon,u_0^{[\leq N]} \circ K_\varepsilon^{[\leq N]}(\theta) - K_\varepsilon^{[\leq N]}(\theta + \omega) = \sum_{n=N+1}^{\infty} \left( \frac{S_{n-1}(\theta)}{S_{n-1}(1)} \right) \varepsilon^n$$

and therefore, for any fixed $\varepsilon$, $E_\varepsilon^{[N,2N]}(\theta)$ is a trigonometric polynomial of degree $2aN$. Moreover, in this case the matrix $M_\varepsilon^{[\leq N]}(\theta) = [D K_\varepsilon^{[\leq N]}(\theta)]^{-1} J K^{[\leq N]}(\theta) D K^{[\leq N]}(\theta) K^{[\leq N]}(\theta)$ is given by

$$M_\varepsilon^{[\leq N]}(\theta) = \begin{bmatrix} 1 + \sum_{k=0}^{N} u_k(\theta) \varepsilon^k & N^{[\leq N]}(\theta) \sum_{k=0}^{N} (u_k(\theta - \omega) - u_k(\theta)) \varepsilon^k \\ \sum_{k=0}^{N} (u_k(\theta) - u_k(\theta - \omega)) \varepsilon^k & N^{[\leq N]}(\theta) \left( 1 + \sum_{k=0}^{N} u_k(\theta) \varepsilon^k \right) \end{bmatrix},$$

where $N^{[\leq N]}(\theta) = \left( 1 + \sum_{k=0}^{N} u_k(\theta) \varepsilon^k \right)^2 + (\sum_{k=0}^{N} (u_k(\theta) - u_k(\theta - \omega)) \varepsilon^k)^2 \right)^{-1}$. So,

$$(M_\varepsilon^{[\leq N]} \circ T_\omega)^{-1} = \begin{bmatrix} (N^{[\leq N]} \circ T_\omega) \left( 1 + \sum_{k=0}^{N} u_k(\theta + \omega) \varepsilon^k \right) & \sum_{k=0}^{N} (u_k(\theta) - u_k(\theta + \omega)) \varepsilon^k \\ \sum_{k=0}^{N} (u_k(\theta) - u_k(\theta + \omega)) \varepsilon^k & 1 + \sum_{k=0}^{N} u_k(\theta + \omega) \varepsilon^k \end{bmatrix}$$

which implies that $E_\varepsilon^{[N,2N]}(\omega)$ is a trigonometric polynomial of degree $3aN$. Remember that $E_\varepsilon^{[N,2N]}$ is the second row of the vector $E_\varepsilon^{[N,2N]} = (M_\varepsilon^{[\leq N]} \circ T_\omega)^{-1} E_\varepsilon^{[N,2N]}$. Note that $J = \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix}$.
Furthermore, we have $D_{\mu} f_{\varepsilon, \mu}^{(N)}(x, y) = \begin{pmatrix} 1 \\ 1 \end{pmatrix}$, then the second row, $\tilde{\Lambda}^N_{\varepsilon, 2}$, of the vector $\tilde{\Lambda}^N_{\varepsilon}$ = $(M_{\varepsilon}^{[\leq N]} \circ T_\varepsilon)^{-1} D_{\mu} f_{\varepsilon, \mu}^{(N)} \circ K^n_{\varepsilon}^{[\leq N]}$ is a trigonometric polynomial of degree $aN$.

The following proposition summarizes the computations presented above and assures that hypothesis HTP1 and HTP2 of theorem B are satisfied for the dissipative standard map.

**Proposition 57.** For any $N \in \mathbb{N}$, if $V(\theta)$ in (A.1) is a trigonometric polynomial of degree $a$, then $E_{\varepsilon, 2}^{(N, 2aN)}$ is a trigonometric polynomial of degree $3aN$, $\tilde{\Lambda}^N_{\varepsilon, 2}$ is a trig polynomial of degree $aN$, and

$$E_{\varepsilon, 2}^{(N, 2aN)}(\theta) \equiv DK^{[\leq N]}_{\varepsilon}(\theta + \omega)^{\top} J \circ K^{[\leq N]}_{\varepsilon}(\theta) DK^{[\leq N]}_{\varepsilon}(\theta + \omega)$$

$$- D(f_{\varepsilon, \mu}^{(N)} \circ K^{[\leq N]}_{\varepsilon}(\theta))^{\top} J \circ (f_{\varepsilon, \mu}^{(N)} \circ K^{[\leq N]}_{\varepsilon}(\theta))$$

$$\times D(f_{\varepsilon, \mu}^{(N)} \circ K^{[\leq N]}_{\varepsilon}(\theta))$$

(A.9)

is a trigonometric polynomial of degree $2aN$.

**Proof.** It is only left to prove the last claim. Note that $E_{\varepsilon, 2}^{(N)}(\theta)$ is the expression in coordinates of $(K^{[\leq N]}_{\varepsilon} \circ T_\varepsilon)^{\top} \Omega - (f_{\varepsilon, \mu}^{(N)} \circ K^{[\leq N]}_{\varepsilon})^{\top} \Omega$. Now, using the fact that $f_{\varepsilon, \mu}$ is conformally symplectic we have $(f_{\varepsilon, \mu}^{(N)} \circ K^{[\leq N]}_{\varepsilon})^{\top} \Omega = K^{[\leq N]}_{\varepsilon}^{\top} f_{\varepsilon, \mu}^{(N)} \circ K^{[\leq N]}_{\varepsilon} \Omega = \lambda(\varepsilon) K^{[\leq N]}_{\varepsilon}^{\top} \Omega$, which means that, in coordinates

$$E_{\varepsilon, 2}^{(N)}(\theta, \varepsilon) = DK^{[\leq N]}_{\varepsilon}(\theta + \omega)^{\top} J \circ K^{[\leq N]}_{\varepsilon}(\theta + \omega) DK^{[\leq N]}_{\varepsilon}(\theta + \omega)$$

$$- \lambda(\varepsilon) DK^{[\leq N]}_{\varepsilon}(\theta)^{\top} J \circ K^{[\leq N]}_{\varepsilon}(\theta) DK^{[\leq N]}_{\varepsilon}(\theta)$$

(A.10)

which is a polynomial of degree $2aN$ due to the fact that $J$ is a constant matrix and

$$DK^{[\leq N]}_{\varepsilon}(\theta) = \begin{pmatrix} 1 + \sum_{n=0}^{N} u_\varepsilon(\theta) \varepsilon^n \\
\sum_{n=0}^{N} (u_\varepsilon(\theta) - u_\varepsilon(\theta - \omega)) \varepsilon^n \end{pmatrix}$$

is a trigonometric polynomial of degree $aN$. \qed

**A.2. Uniqueness**

Note that for $\varepsilon = 0$, $M_0 = I$. Also note that the coefficients of the expansion (A.8) are given by

$$K_\mu(\theta) = \begin{pmatrix} u_\mu(\theta) \\
\mu_\mu(\theta) - u_\mu(\theta - \omega) \end{pmatrix} \quad \text{for } n \geq 1.$$

Therefore, the normalization condition

$$\int_T [M_0^{-1} K_\mu(\theta)]_1 d\theta = 0$$
in this case has the form
\[
\int_T u_\epsilon(\theta) d\theta = 0,
\]
which is satisfied by the construction of the \( u_\epsilon \)'s. Thus, the expansion given in (A.8) is the only one which satisfies the normalization condition.

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