# KAM Lectures* 

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The aim of these lectures is to present, in a self contained way, the fundamentals of KAM theory, which, as well known, deals with the problem of constructing quasi-periodic motions in real-analytic or smooth conservative dynamical systems.

KAM theory is based upon quantitative techniques designed to overcome the socalled small denominator difficulties arising in the construction of quasi-periodic motions and works under rather stringent smallness and regularity assumptions.

For sake of presentation, we will consider only second order Hamiltonian systems with a finite number of degrees of freedom (periodic in the "space" variables), i.e., systems governed by Hamiltonian functions of of the form

$$
\begin{equation*}
H(y, x)=\frac{y^{2}}{2}+V(x) \tag{1}
\end{equation*}
$$

where $y$ and $x$ are standard symplectic variables $(y, x) \in \mathbb{R}^{d} \times \mathbb{T}^{d}$, and $V: \mathbb{T}^{d} \rightarrow \mathbb{R}$ is a (multi-periodic) smooth or real-analytic function; $y^{2}:=y \cdot y:=\sum_{j=1}^{d} y_{j}^{2}$. Here, $\mathbb{T}^{d}$ denotes the standard flat $d$-torus $\mathbb{T}^{d}:=\mathbb{R}^{d} /\left(2 \pi \mathbb{Z}^{d}\right)$; the (standard) symplectic structure is: $d y \wedge d x=\sum_{j=1}^{d} d y_{j} \wedge d x_{j}$ and the Hamilton equations are

$$
\begin{equation*}
\dot{y}=-H_{x}, \quad \dot{x}=H_{y}, \tag{2}
\end{equation*}
$$

[^0]where $H_{y}$ denotes the $y$-gradient $\left(H_{y_{1}}, \ldots, H_{y_{d}}\right)$ and $H_{x}$ denotes the $x$-gradient ( $H_{x_{1}}, \ldots, H_{x_{d}}$ ); dot denotes time derivative.

The point of view taken up in these lectures is that of non-linear functional analysis, as we briefly proceed to explain. The problem of constructing (maximal ${ }^{1}$ ) quasiperiodic solutions is essentially equivalent to solve a non-linear partial differential equation on $\mathbb{T}^{d}, \mathcal{E}(u)=0$, with real-analytic or $C^{k}$ coefficients. If one is given an approximate solution, i.e. a function $v$ for which $\mathcal{E}(v)$ is not zero but small (in suitable norms), then, under suitable conditions, it is possible to find a true near-by solution. The method we shall follow is based on a Newton ("quadratic") scheme, which allow to construct a sequence of better and better approximations (living in larger and larger Banach spaces) converging to a true solution. The loss of regularity (related to the inversion of non elliptic differential operators and to the above mentioned small denominator problems) arising in solving the associated linearized equation is overcome by the speed of convergence of the scheme.

The approach presented here - sometimes referred to as KAM theory in configuration space - avoids completely the use of symplectic transformations and needs less preparation than standard KAM theory.

The notes of the lectures are divided in two chapters:
In the first chapter a KAM theorem establishing the existence of quasi-periodic solutions (with prescribed "diophantine" frequencies), in real-analytic setting, is presented. The "potential" $V$ in (1) is not assumed to be small; what allows to start up the perturbative procedure is the existence of a good enough approximate solution.
While no effort is put in trying to get "optimal estimates", a certain care is devoted to perform explicit estimates and also to discuss convenient norms (Fourier and complex sup-norms).

In the second chapter, we shall consider Hamiltonians $H$ in (1) with $V \in C^{l}\left(\mathbb{T}^{d}\right)$, which shall be assumed to be small in $C^{1}$ norm. Then, assuming $l$ big enough and using the approximation technique due to Bernstein, Jackson, Moser and Zehnder ${ }^{2}$, we shall construct (using the real-analytic KAM theorem of the first chapter) a sequence of real-analytic approximate solutions converging to $C^{s}$ quasi-periodic solutions; explicit estimates on $l$ and $s$ will be given.

[^1]The main references are:
[1] D. Salamon, E. Zehnder : KAM theory in configuration space, Comm. Math. Helv. 64 (1989), 84-132
[2] D. Salamon, The Kolmogorov-Arnold-Moser theorem, FIM-Preprint, ETHZurich, (1986), available on
http://www.math.ethz.ch/~salamon/PREPRINTS/KAM.htm
For the analytic part, see also:
[3] A. Celletti and L. Chierchia: A constructive theory of Lagrangian tori and computer-assisted applications, Dynamics reported, 60-130, Dynam. Report. Expositions Dynam. Systems (N.S.), 4, Springer, Berlin, 1995

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## 1 Analytic KAM Theory

### 1.1 Warm up: Newton scheme for the standard IFT

The aim of this section is to discuss a proof of the (standard) Implicit Function Theorem in $\mathbb{R}^{n}$ based on the "Newton method" with the purpose of illustrating, in a trivial case, the scheme of proof that we shall use to construct quasi-periodic motions for Hamiltonian systems.

Let $(\bar{y}, \bar{x}) \in \mathbb{R}^{n} \times \mathbb{R}^{m}$; denote by $D_{\rho}^{n}$ the closed ball in $\mathbb{R}^{n}$ centered at $\bar{y}$ with radius $\rho$ and by $D_{\bar{r}}^{m}$ the closed ball in $\mathbb{R}^{m}$ centered at $\bar{x}$ with radius $\bar{r}$; let
$X_{r, \rho}$ denote the Banach space, $C\left(D_{r}^{m}, D_{\rho}^{n}\right)$, of continuous function from $D_{r}^{m}$ into $D_{\rho}^{n}$ endowed with the sup-norm.

Theorem 1 Let $F \in C\left(D_{\rho}^{n} \times D_{\bar{r}}^{m}, \mathbb{R}^{n}\right)$ be such that $y \mapsto F(y, x) \in C^{2}\left(D_{\rho}^{n}, \mathbb{R}^{n}\right)$ for all $x \in D_{\bar{r}}^{m}$ with $F_{y}$ invertible on $D_{\rho}^{n} \times D_{\bar{r}}^{m}$. Let $\alpha$ and $\beta$ be positive numbers such that

$$
\begin{equation*}
\left\|\left(F_{y}\right)^{-1}\right\|_{\rho, \bar{r}} \leq \alpha, \quad \frac{\alpha^{2}}{2}\left\|F_{y y}\right\|_{\rho, \bar{r}} \leq \beta, \tag{3}
\end{equation*}
$$

$\|\cdot\|_{\rho, \bar{r}}$ being short for $\sup _{D_{\rho}^{n} \times D_{\bar{T}}^{m}}|\cdot|$. Suppose that, for some $0<r \leq \bar{r}$ and $0<\sigma<1$, there exists $u_{0} \in X_{r, \rho}$ such that $\left\|u_{0}-\bar{y}\right\|_{r}:=\sup _{D_{r}^{m}}\left|u_{0}-\bar{y}\right|<\rho$ and:

$$
\begin{equation*}
\left\|F\left(u_{0}(x), x\right)\right\|_{r} \leq \min \left\{\frac{\sigma}{\beta}, \frac{1-\sigma}{\alpha}\left(\rho-\left\|u_{0}-\bar{y}\right\|_{r}\right)\right\} . \tag{4}
\end{equation*}
$$

Then, there exists a unique function $u \in X_{r, \rho}$ such that:

$$
\begin{equation*}
F(u(x), x)=0, \quad \forall x \in D_{r}^{m}, \tag{5}
\end{equation*}
$$

and

$$
\left\|u-u_{0}\right\|_{r} \leq \frac{\alpha}{1-\sigma}\left\|F\left(u_{0}(x), x\right)\right\|_{r}
$$

Remark 1 (i) The limiting case $\beta=0$ corresponds to the linear case

$$
F(y, x)=a(x)+A(x) y
$$

(with $A$ invertible), in which case the solution of $F(u, x)=0$ is simply $u=-A^{-1} a$.
(ii) If $F(\bar{y}, \bar{x})=0$, one can obviously take $u_{0}(x) \equiv \bar{y}$ (choosing suitably $r$ so as to meet condition (4)).
(iii) The function $u_{0}$ is called an approximate solution of (5); the function

$$
\begin{equation*}
\varepsilon_{0}(x):=F\left(u_{0}(x), x\right), \tag{6}
\end{equation*}
$$

is the associated error function. The inequality (4) should be interpreted as a smallness condition on the error function and the IFT can be rephrased by saying that if the smallness condition is verified by the error function $\varepsilon_{0}$ associated to the approximate solution $u_{0}$, than there exists a (unique) true solution $u$, which is $\left\|\varepsilon_{0}\right\|$-close to the approximate solution $u_{0}$.

Proof We first show how to construct out of $u_{0}$ a new approximate solution $u_{1}$ for which the associated error function $\varepsilon_{1}(x):=F\left(u_{1}(x), x\right)$ is quadratically smaller that $\varepsilon_{0}$.

Let $\varepsilon_{0}$ be as in (6) and define

$$
\begin{equation*}
w_{0}(x)=-\left(F_{y}\left(u_{0}(x), x\right)\right)^{-1} \varepsilon_{0}(x), \quad u_{1}:=u_{0}+w_{0} . \tag{7}
\end{equation*}
$$

We claim that $u_{1} \in X_{r, \rho}$ and that $w_{0}$ and $\varepsilon_{1}:=F\left(u_{1}, x\right)$ verify:

$$
\begin{equation*}
\left\|w_{0}\right\|_{r} \leq \alpha\left\|\varepsilon_{0}\right\|_{r}, \quad\left\|\varepsilon_{1}\right\|_{r} \leq \beta\left\|\varepsilon_{0}\right\|_{r}^{2} \tag{8}
\end{equation*}
$$

In fact, the first estimate in (8) is immediate consequence of the definitions of $w_{0}$ and $\alpha$. To show that $u_{1}: D_{r}^{m} \rightarrow D_{\rho}^{n}$, we compute:

$$
\begin{aligned}
\left\|u_{1}-\bar{y}\right\|_{r} & :=\left\|u_{0}+w_{0}-\bar{y}\right\|_{r} \leq\left\|u_{0}-\bar{y}\right\|_{r}+\left\|w_{0}\right\|_{r} \\
& \leq\left\|u_{0}-\bar{y}\right\|_{r}+\alpha\left\|\varepsilon_{0}\right\|_{r} \leq\left\|u_{0}-\bar{y}\right\|_{r}+\alpha \frac{1-\sigma}{\alpha}\left(\rho-\left\|u_{0}-\bar{y}\right\|_{r}\right) \\
& <\rho,
\end{aligned}
$$

where we have used the assumption (4) on $\varepsilon_{0}:=F\left(u_{0}, x\right)$. Observe that, by the definition of $\varepsilon_{1}, w_{0}$ and Taylor's formula, one gets:

$$
\begin{align*}
\varepsilon_{1} & :=F\left(u_{1}, x\right):=F\left(u_{0}+w_{0}, x\right)=F\left(u_{0}, x\right)+F_{y}\left(u_{0}, x\right) w_{0}+Q \\
& =\varepsilon_{0}+F_{y}\left(u_{0}, x\right) w_{0}+Q=Q \tag{9}
\end{align*}
$$

where ${ }^{3}$

$$
\begin{equation*}
Q=\int_{0}^{1}(1-t) F_{y y}\left(u_{0}+t w_{0}, x\right) w_{0} \cdot w_{0} d t \tag{10}
\end{equation*}
$$

Thus, by the estimates on $w_{0}$ in (8) and the definition of $\beta$, we get

$$
\begin{equation*}
\left\|\varepsilon_{1}\right\|_{r}=\|Q\|_{r} \leq \frac{1}{2}\left\|F_{y y}\right\|_{\rho, \bar{r}}\left\|w_{0}\right\|_{r}^{2} \leq \frac{\alpha^{2}}{2}\left\|F_{y y}\right\|_{\rho, \bar{r}}\left\|\varepsilon_{0}\right\|_{r}^{2} \leq \beta\left\|w_{0}\right\|_{r}^{2} \tag{11}
\end{equation*}
$$

completing the proof of (8).
The idea is, now, to iterate such construction: Fix $k \geq 2$ and assume that $u_{1}, \ldots, u_{k-1}$ are given approximate solutions belonging to the Banach space $X_{r, \rho}$ and such that, if one defines

$$
\begin{equation*}
w_{j}:=u_{j+1}-u_{j}, \quad \varepsilon_{j}(x):=F\left(u_{j}(x), x\right), \quad(0 \leq j \leq k-2) \tag{12}
\end{equation*}
$$

[^2]then the following inequalities hold for all $0 \leq j \leq k-2$ :
\[

$$
\begin{equation*}
\left\|w_{j}\right\|_{r} \leq \alpha\left\|\varepsilon_{j}\right\|_{r}, \quad\left\|\varepsilon_{j+1}\right\|_{r} \leq \beta\left\|\varepsilon_{j}\right\|_{r}^{2} \tag{13}
\end{equation*}
$$

\]

Note that such inductive assumption has been verified for $k=2$ with $u_{1}$ as in (7).

We claim that, under the inductive assumption (12) and (13), setting

$$
\begin{equation*}
\varepsilon_{k-1}(x):=F\left(u_{k-1}(x), x\right), \quad w_{k-1}(x)=-F_{y}\left(u_{k-1}(x), x\right) \varepsilon_{k-1}(x) \tag{14}
\end{equation*}
$$

then one has

$$
\begin{equation*}
u_{k}:=u_{k-1}+w_{k-1} \in X_{r, \rho} \tag{15}
\end{equation*}
$$

and (13) holds also for $j=k-1$.
In fact, the estimate on $\left\|w_{k-1}\right\|_{r}$ follows at once (as above) from the definition of $w_{k-1}$ and $\alpha$ (and the inductive assumption on $u_{k-1}$ ). Let us, now, show (15). Multiplying by $\beta$ the second relation in (13) can be rewritten as

$$
\begin{equation*}
\beta\left\|\varepsilon_{j+1}\right\|_{r} \leq\left(\beta\left\|\varepsilon_{j}\right\|_{r}\right)^{2} \tag{16}
\end{equation*}
$$

which iterated leads to

$$
\begin{equation*}
\beta\left\|\varepsilon_{j}\right\|_{r} \leq\left(\beta\left\|\varepsilon_{0}\right\|_{r}\right)^{2^{j}}, \quad \forall 0 \leq j \leq k-1 \tag{17}
\end{equation*}
$$

Thus, by (12), (13) (first inequality), (17) and (4), one has

$$
\begin{aligned}
\left\|u_{k}-\bar{y}\right\|_{r} & =\left\|u_{0}+\sum_{j=0}^{k-1} w_{j}-\bar{y}\right\|_{r} \\
& \leq\left\|u_{0}-\bar{y}\right\|_{r}+\sum_{j=0}^{k-1}\left\|w_{j}\right\|_{r} \\
& \leq\left\|u_{0}-\bar{y}\right\|_{r}+\alpha \sum_{j=0}^{k-1}\left\|\varepsilon_{j}\right\|_{r} \\
& \leq\left\|u_{0}-\bar{y}\right\|_{r}+\frac{\alpha}{\beta} \sum_{j=0}^{k-1}\left(\beta\left\|\varepsilon_{0}\right\|_{r}\right)^{2^{j}} \\
& \leq\left\|u_{0}-\bar{y}\right\|_{r}+\frac{\alpha}{\beta} \sum_{j=1}^{\infty}\left(\beta\left\|\varepsilon_{0}\right\|_{r}\right)^{j}
\end{aligned}
$$

$$
\begin{aligned}
& =\left\|u_{0}-\bar{y}\right\|_{r}+\alpha \frac{\left\|\varepsilon_{0}\right\|_{r}}{1-\beta\left\|\varepsilon_{0}\right\|_{r}} \\
& \leq\left\|u_{0}-\bar{y}\right\|_{r}+\alpha \frac{\left\|\varepsilon_{0}\right\|_{r}}{1-\sigma} \\
& \leq\left\|u_{0}-\bar{y}\right\|_{r}+\alpha \frac{1-\sigma}{\alpha} \frac{\rho-\left\|u_{0}-\bar{y}\right\|_{r}}{1-\sigma}=\rho .
\end{aligned}
$$

This shows (15). At this point, also the estimate on $\left\|\varepsilon_{k}\right\|_{r}$ follows: just replace $\varepsilon_{1}, u_{0}$ and $w_{0}$ in (9) $\div(11)$ by, respectively, $\varepsilon_{k}, u_{k-1}$ and $w_{k-1}$.

Thus, thanks to (4), the construction can be iterated indefinitely and $\left\{u_{k}\right\}$ will converge to a function $u \in X_{r, \rho}$. Clearly, since $\left\|\varepsilon_{j}\right\|_{r} \rightarrow 0$ (superexponentially fast), one has

$$
F(u, x)=\lim F\left(u_{k}, x\right)=\lim \varepsilon_{k}=0,
$$

showing (5).
Uniqueness is an obvious consequence of the invertibility of $F_{y}$.

Remark 2 (i) The approximate solutions $u_{k}$ 's belong to the same Banach space $X_{r, \rho}$. This is so because $w_{k-1}$ belongs to the same space of $u_{k-1}$. In the more complicate case of quasi-periodic solutions for Hamiltonian systems this will not be the case any more: the analogous of $F_{y}^{-1}$ will be an unbounded operator (involving small divisors) and (the analogous of) $w_{k-1}$ will lie, in general, in a smaller Banach space.
(ii) In fact, even formally, it will not be possible to solve the linearized equation ${ }^{4}$ exactly but only up to quadratically small terms.
(iii) The argument to prove (local) uniqueness in the quasi-periodic case will be different (because of the lack of invertibility of $F_{y}$ ).

### 1.2 Quasi-periodic solutions (definitions)

Let, as above, $\mathbb{T}^{d}:=\mathbb{R}^{d} /\left(2 \pi \mathbb{Z}^{d}\right)$ be the standard $d$-dimensional flat torus and let $\Omega$ be a bounded domain in $\mathbb{R}^{d}$. Consider a smooth (say $C^{2}$ ) Hamiltonian

[^3]$H(y, x)$ from $\Omega \times \mathbb{T}^{d}$ to $\mathbb{R}$ and the associated Hamiltonian equations
\[

\left\{$$
\begin{array}{l}
\dot{y}_{i}=-\frac{\partial H}{\partial x_{i}}  \tag{18}\\
\dot{x}_{i}=\frac{\partial H}{\partial y_{i}}
\end{array}
$$ \quad i=1, ···, d .\right.
\]

An interesting example is when the system is nearly-integrable, i.e., when $H$ is of the form

$$
H(y, x)=H_{0}(y)+\varepsilon H_{1}(y, x)
$$

with $\varepsilon$ a small parameter. The corresponding Hamiltonian equations become

$$
\left\{\begin{array}{l}
\dot{y}_{i}=-\varepsilon \frac{\partial H_{1}}{\partial x_{i}}  \tag{19}\\
\dot{x}_{i}=\frac{\partial H_{0}}{\partial y_{i}}+\varepsilon \frac{\partial H_{1}}{\partial y_{i}}
\end{array} \quad i=1, \ldots, d .\right.
$$

When $\varepsilon=0$ this system is completely integrable and all solutions,

$$
\left\{\begin{array}{l}
y(t)=y(0), \\
x(t)=x(0)+\frac{\partial H_{0}}{\partial y}(y(0)) t, \quad(\bmod (2 \pi, \ldots, 2 \pi)),
\end{array}\right.
$$

are quasi-periodic:

Definition $1 A$ solution $(y(t), x(t))$ of (18) is said to be quasi-periodic, if there exist a vector $\omega \in \mathbb{R}^{d}$ (frequency vector) and two functions $u, v \in$ $C^{2}\left(\mathbb{T}^{d}, \mathbb{R}^{d}\right)$ such that

$$
\left\{\begin{array}{l}
y(t)=v(\omega t)  \tag{20}\\
x(t)=\omega t+u(\omega t) \quad(\bmod (2 \pi, \ldots, 2 \pi)) .
\end{array}\right.
$$

for every $t$. If the frequency vector is rationally independent (i.e. ${ }^{5} \omega \cdot n \neq 0$ for every $n \in \mathbb{Z}^{d} \backslash\{0\}$ ), then the solution $(y(t), x(t))$ is said to be maximal quasi-periodic.

Remark 3 (i) Non-maximal quasi-periodic solutions include periodic solutions: this is the case when there exist $T>0$ and $n \in \mathbb{Z}^{d}$ such that $\omega T=2 \pi n$; notice that in this case there exist $d-1$ linearly independent vectors $n_{j} \in \mathbb{Z}^{d}$ such that $\omega \cdot n_{j}=0$ for $^{6} j=1, \ldots, d-1$. More in general, frequencies may be

[^4]classified in terms of the number of relations $\omega \cdot m=0$ satisfied by $\omega$ with independent vectors $m \in \mathbb{Z}^{d}$.
(ii) A maximal quasi-periodic solution is said to be non-degenerate if ( $\omega$ is rationally independent and) the map $\theta \in \mathbb{T}^{d} \rightarrow \theta+u(\theta) \in \mathbb{T}^{d}$ is a diffeomorphism of $\mathbb{T}^{d}$ so that the $\operatorname{map} \theta \in \mathbb{T}^{d} \rightarrow(v(\theta), \theta+u(\theta)) \in \mathbb{R}^{d} \times \mathbb{T}^{d}$ yields an embedding of the $d$-dimensional torus into the phase space $\Omega \times \mathbb{T}^{d}$. The relation (20) says that non-degenerate maximal quasi-periodic solutions correspond to $d$-dimensional invariant tori on which the $H$-flow is conjugate to the linear flow $\theta \rightarrow \theta+\omega t$.
(iii) In these lectures we shall consider only non-degenerate maximal quasiperiodic solutions and, hereafter, "quasi-periodic solution" will be used as synonymous of "non-degenerate maximal quasi-periodic solutions". In particular, the frequency vector $\omega$ is always assumed to be rationally independent.

Consider a quasi-periodic solution $(y(t), x(t))$ as in (20). Differentiating it with respect to $t$ we get

$$
\left\{\begin{array}{l}
\dot{y}(t)=D v(\omega t) \\
\dot{x}(t)=\omega+D u(\omega t)
\end{array}\right.
$$

where

$$
D:=D_{\omega}:=\sum_{i=1}^{d} \omega_{i} \frac{\partial}{\partial \theta_{i}}
$$

Since $(y, x)$ is a solution of (18), we have

$$
\left\{\begin{array}{l}
D v_{i}(\omega t)=-\frac{\partial H}{\partial x_{i}}(v(\omega t), \omega t+u(\omega t)) \\
\omega_{i}+D u_{i}(\omega t)=\frac{\partial H}{\partial y_{i}}(v(\omega t), \omega t+u(\omega t))
\end{array} \quad i=1, \ldots, d\right.
$$

which, by density of the trajectory $t \mapsto \omega t$ on ${ }^{7} \mathbb{T}^{d}$, are equivalent to

$$
\left\{\begin{array}{l}
D v(\theta)=-H_{x}(v(\theta), \theta+u(\theta))  \tag{21}\\
\omega+D u(\theta)=H_{y}(v(\theta), \theta+u(\theta))
\end{array} \quad \theta \in \mathbb{T}^{d}\right.
$$

[^5]Hereafter, (for simplicity) we shall consider only Hamiltonians $H$ of the form

$$
H(y, x)=\frac{y^{2}}{2}+V(x):=\frac{1}{2} \sum_{j=1}^{d} y_{j}^{2}+V(x) .
$$

In this special case (18) takes the form

$$
\left\{\begin{array}{l}
\dot{y}=-V_{x} \\
\dot{x}=y
\end{array}\right.
$$

or, equivalently,

$$
\ddot{x}=-V_{x} .
$$

Notice that, in such a case, the second equation in (21) becomes simply

$$
\begin{equation*}
v(\theta)=\omega+D u(\theta) \tag{22}
\end{equation*}
$$

so that the system (21) becomes the following single (vector) equation for $u$ :

$$
\begin{equation*}
D^{2} u(\theta)=-V_{x}(\theta+u(\theta)) . \tag{23}
\end{equation*}
$$

The lectures are devoted to discuss solutions of (23).
It is clear that an important rôle in the study of (21) or (23) is played by the linear equation

$$
D u=f
$$

with $f$ a given function on $\mathbb{T}^{d}$. Proceeding formally, we expand both sides in Fourier series getting

$$
\sum_{n \in \mathbb{Z}^{d}} f_{n} e^{i n \cdot \theta}=\sum_{n \in \mathbb{Z}^{d}} i(\omega \cdot n) u_{n} e^{i n \cdot \theta}
$$

Equating Fourier coefficient, we get, for $n=0$, the compatibility condition ${ }^{8}$

$$
\begin{equation*}
f_{0}=\langle f\rangle=0 \tag{24}
\end{equation*}
$$

and, for $n \neq 0$,

$$
\begin{equation*}
u_{n}=\frac{f_{n}}{i(\omega \cdot n)}, \quad\left(n \in \mathbb{Z}^{d} \backslash\{0\}\right) \tag{25}
\end{equation*}
$$

${ }^{8}$ We denote $\langle\cdot\rangle:=f_{\mathbb{T}^{d}} \cdot d \theta:=(2 \pi)^{-d} \int_{\mathbb{T}^{d}} \cdot d \theta$.

The denominator $(\omega \cdot n)$ in (25), even though never vanishes, might become arbitrarily small making doubtful the convergence of the Fourier series

$$
\begin{equation*}
\sum_{n \in \mathbb{Z}^{d}} u_{n} e^{i n \cdot \theta} . \tag{26}
\end{equation*}
$$

Definition 2 We say that $\omega \in \mathbb{R}^{d}$ is $(\gamma, \tau)$-diophantine if $\gamma, \tau$ are positive constants such that

$$
\begin{equation*}
|\omega \cdot n| \geq \frac{\gamma}{|n|^{\tau}} \quad \text { for every } \quad n \in \mathbb{Z}^{d} \backslash\{0\} \tag{27}
\end{equation*}
$$

Remark 4 For $\tau>d-1$ fixed, the set of diophantine vectors is of full measure (exercise). For $\tau<d-1$ (27) is never satisfied (Liouville).

Suppose now that $\omega$ is $(\gamma, \tau)$-diophantine and $f$ is a smooth enough function with vanishing mean value, $\langle f\rangle=0$. Then (26)-(25) actually define the function $u$, solution of $D u=f$, up to an additive constant (the average of $u$ ); the unique solution of the system:

$$
D u=f, \quad\langle u\rangle=0,
$$

will be denoted by $D^{-1} f$.
Exercise Find a lower bound on $k$ so that if $f \in C^{k}\left(\mathbb{T}^{d}\right)$ then $D^{-1} f$ has an absolutely convergent Fourier series expansion.

Remark 5 The analysis described in these lectures could be easily extended to the non-autonomous case, i.e, the case when the potential $V=V(x, t)$ depends also explicitly (and periodically) on time $t, V: \mathbb{T}^{d+1} \rightarrow \mathbb{R}$. In such a case $D_{\omega}$ has to be replaced by

$$
D=\sum_{i=1}^{d} \omega_{i} \frac{\partial}{\partial x_{i}}+\frac{\partial}{\partial t}
$$

with $(\omega, 1) \in \mathbb{R}^{d+1}$ rationally independent and equation (23) becomes

$$
D^{2} u(\theta, t)=-V_{x}(\theta+u(\theta, t), t) .
$$

### 1.3 Newton scheme for Quasi-periodic solutions

In this section we describe the Newton scheme on which the construction of solutions of the functional equation (23) will be based.

The strategy that we shall follow mimics the proof of Theorem 1: we shall start from an approximate solution $v$ of (23), i.e., a (smooth) function $v$ such that the associated error function

$$
\begin{equation*}
\varepsilon:=\mathcal{E}(v):=D^{2} v+V_{x}(\theta+v), \tag{28}
\end{equation*}
$$

is "small" and try to construct a "better" approximate solution

$$
\begin{equation*}
v^{\prime}:=v+w, \tag{29}
\end{equation*}
$$

whose associated error function

$$
\begin{equation*}
\varepsilon^{\prime}:=\mathcal{E}\left(v^{\prime}\right):=D^{2} v^{\prime}+V_{x}\left(\theta+v^{\prime}\right) \tag{30}
\end{equation*}
$$

is "quadratically smaller" than the error function associated to $v$.
Remark 6 The discussion in this section will be algebraic in character and the necessary estimates will be discussed later (§1.5). Therefore, words such as "small" or "quadratically smaller" are used, here, in a somewhat formal way ${ }^{9}$. Roughly speaking, the idea is to look for $w \sim \varepsilon$ (i.e., "of the same order of $\varepsilon^{\prime \prime}$ ) so that $\mathcal{E}(v+w) \sim \varepsilon^{2}$. However, as clarified also in Remark 7 below, the reader can also disregard any reference to "smallness" following only the algebraic identities involved.

Define $Q_{1}$ as

$$
\begin{equation*}
Q_{1}:=V_{x}(\theta+v+w)-V_{x}(\theta+v)-V_{x x}(\theta+v) w . \tag{31}
\end{equation*}
$$

and note that, by Taylor's formula, $Q_{1}$ is quadratic in ${ }^{10} w$. Expanding $V_{x}(\theta+$ $v+w)$ we find:

$$
\varepsilon^{\prime}:=\mathcal{E}\left(v^{\prime}\right):=D^{2} v+D^{2} w+V_{x}(\theta+v+w)
$$

[^6]\[

$$
\begin{align*}
& =D^{2} v+V_{x}(\theta+v)+D^{2} w+V_{x x}(\theta+v) w+Q_{1} \\
& =: \\
& =: \quad \mathcal{E}(v)+D^{2} w+V_{x x}(\theta+v) w+Q_{1}  \tag{32}\\
& =: D^{2} w+V_{x x}(\theta+v) w+Q_{1} .
\end{align*}
$$
\]

The perfect analogue of the Newton scheme described in the proof of the standard IFT given in § 1 would consist in finding an "explicit" solution of the the following PDE on $\mathbb{T}^{d}$

$$
\begin{equation*}
\varepsilon+D^{2} w+V_{x x}(\theta+v) w=0 \tag{33}
\end{equation*}
$$

However this is not so easy and, in fact, we shall be able to solve (33) only up to quadratic terms in $\varepsilon$.

To proceed further, we look at the variation equation for (23), i.e., the equation

$$
\begin{equation*}
\varepsilon_{\theta}=D^{2} v_{\theta}+V_{x x}(\theta+v)\left(I+v_{\theta}\right), \tag{34}
\end{equation*}
$$

which is gotten by differentiating with respect to $\theta$ the system (23); here, for a given function $u: \mathbb{T}^{d} \rightarrow \mathbb{R}^{d}$, $u_{\theta}$ denotes the Jacobian matrix

$$
u_{\theta}=\left(\frac{\partial u_{i}}{\partial \theta_{j}}(\theta)\right)_{i, j=1, \ldots, d}
$$

and $I:=I_{d}$ denotes the unit $(d \times d)$ matrix.
Setting

$$
\begin{equation*}
M:=I+v_{\theta} \tag{35}
\end{equation*}
$$

we can rewrite (34) in the form

$$
\begin{equation*}
\varepsilon_{\theta}=D^{2} M+V_{x x}(\theta+v) M \tag{36}
\end{equation*}
$$

Assume that $M(\theta)$ is invertible for all $\theta \in \mathbb{T}^{d}$. From (36) we get

$$
V_{x x}(\theta+v)=\left(\varepsilon_{\theta}-D^{2} M\right) M^{-1}
$$

and plugging this equality in (32), we find

$$
\begin{align*}
\varepsilon^{\prime} & =\varepsilon+D^{2} w+\left(\varepsilon_{\theta}-D^{2} M\right) M^{-1} w+Q_{1} \\
& =\varepsilon+D^{2} w-\left(D^{2} M\right) M^{-1} w+\varepsilon_{\theta} M^{-1} w+Q_{1} \\
& =: \varepsilon+D^{2} w-\left(D^{2} M\right) M^{-1} w+Q_{2} \tag{37}
\end{align*}
$$

with

$$
\begin{equation*}
Q_{2}=Q_{1}+\varepsilon_{\theta} M^{-1} w \tag{38}
\end{equation*}
$$

Setting

$$
\begin{equation*}
z:=M^{-1} w \tag{39}
\end{equation*}
$$

we get:

$$
\begin{align*}
\varepsilon^{\prime} & =\varepsilon+D^{2}(M z)-\left(D^{2} M\right) z+Q_{2} \\
& =\varepsilon+D(M D z)+D(D M z)-\left(D^{2} M\right) z+Q_{2} \\
& =\varepsilon+D(M D z)+(D M)(D z)+Q_{2} \tag{40}
\end{align*}
$$

Denote by $M^{T}$ the transpose of the matrix $M$ and let $M^{-T}:=\left(M^{T}\right)^{-1}$. Then:

$$
\begin{align*}
\varepsilon^{\prime} & =M^{-T}\left(M^{T} \varepsilon+M^{T} D(M D z)+M^{T}(D M)(D z)\right)+Q_{2} \\
& =M^{-T}\left(M^{T} \varepsilon+D\left(M^{T} M D z\right)-\left(D M^{T}\right)(M D z)+M^{T}(D M)(D z)\right)+Q_{2} \\
& =: M^{-T}\left(M^{T} \varepsilon+D\left(M^{T} M D z\right)\right)+g+Q_{2}, \tag{41}
\end{align*}
$$

with

$$
\begin{equation*}
g:=M^{-T}\left(M^{T} D M-\left(D M^{T}\right) M\right) D z \tag{42}
\end{equation*}
$$

We claim that $g$ is quadratic in $\varepsilon$. To check this, we, first, remark that

$$
\begin{align*}
\left\langle M^{T} D M-\left(D M^{T}\right) M\right\rangle & =\left\langle M^{T} D v_{\theta}-\left(D v_{\theta}^{T}\right) M\right\rangle \\
& =\left\langle\left(I+v_{\theta}^{T}\right) D v_{\theta}-\left(D v_{\theta}^{T}\right)\left(I+v_{\theta}\right)\right\rangle \\
& =\left\langle v_{\theta}^{T} D v_{\theta}-D v_{\theta}^{T} v_{\theta}\right\rangle \tag{43}
\end{align*}
$$

(since $\langle D u\rangle=0$ for any periodic function $u$ ). Integrating by parts,

$$
\begin{aligned}
\left\langle v_{\theta}^{T} D v_{\theta}-D v_{\theta}^{T} v_{\theta}\right\rangle_{i j} & =\sum_{k=1}^{d} f_{\mathbb{T}^{d}}\left(\frac{\partial v_{k}}{\partial \theta_{i}}\left(D \frac{\partial v_{k}}{\partial \theta_{j}}\right)-\left(D \frac{\partial v_{k}}{\partial \theta_{i}}\right) \frac{\partial v_{k}}{\partial \theta_{j}}\right) \\
& =-\sum_{k=1}^{d} f_{\mathbb{T}^{d}}\left(\frac{\partial^{2} v_{k}}{\partial \theta_{j} \partial \theta_{i}}\left(D v_{k}\right)-\left(D v_{k}\right) \frac{\partial^{2} v_{k}}{\partial \theta_{i} \partial \theta_{j}}\right) \\
& =0
\end{aligned}
$$

showing that

$$
\begin{equation*}
\left\langle M^{T} D M-\left(D M^{T}\right) M\right\rangle=0 . \tag{44}
\end{equation*}
$$

Thus, we can write:

$$
M^{T} D M-\left(D M^{T}\right) M=D^{-1}\left[D\left(M^{T} D M-\left(D M^{T}\right) M\right)\right]
$$

But, by (36),

$$
\begin{aligned}
D\left(M^{T} D M-\left(D M^{T}\right) M\right) & =M^{T} D^{2} M-\left(D^{2} M^{T}\right) M \\
& =-M^{T} V_{x x} M+M^{T} \varepsilon_{\theta}+M^{T} V_{x x} M-\varepsilon_{\theta}^{T} M \\
& =M^{T} \varepsilon_{\theta}-\varepsilon_{\theta}^{T} M
\end{aligned}
$$

showing that

$$
\begin{equation*}
\left\langle M^{T} \varepsilon_{\theta}-\varepsilon_{\theta}^{T} M\right\rangle=0, \tag{45}
\end{equation*}
$$

and that

$$
M^{T} D M-\left(D M^{T}\right) M=D^{-1}\left(M^{T} \varepsilon_{\theta}-\varepsilon_{\theta}^{T} M\right)
$$

Thus

$$
\begin{equation*}
g=M^{-T}\left(D^{-1}\left(M^{T} \varepsilon_{\theta}-\varepsilon_{\theta}^{T} M\right)\right) D z \tag{46}
\end{equation*}
$$

is quadratic in $\varepsilon$. Furthermore (41) can be rewritten as

$$
\begin{equation*}
\varepsilon^{\prime}=M^{-T}\left(M^{T} \varepsilon+D\left(M^{T} M D z\right)\right)+Q_{3} \tag{47}
\end{equation*}
$$

with

$$
\begin{equation*}
Q_{3}:=Q_{2}+M^{-T}\left(D^{-1}\left(M^{T} \varepsilon_{\theta}-\varepsilon_{\theta}^{T} M\right)\right) D z \tag{48}
\end{equation*}
$$

We can now show that the equation

$$
\begin{equation*}
M^{T} \varepsilon+D\left(M^{T} M D z\right)=0 \tag{49}
\end{equation*}
$$

can be explicitly solved.
We have already studied the inversion of the differential operator $D$ and therefore we know that a necessary condition to solve our equation is that the average $\left\langle M^{T} \varepsilon\right\rangle$ of $M^{T} \varepsilon$ over $\mathbb{T}^{d}$ is equal to 0 . This is indeed the case, as we proceed to show.

First, by the definitions of $M$ and $\varepsilon$,

$$
\begin{align*}
\left\langle M^{T} \varepsilon\right\rangle & =\left\langle\left(I+v_{\theta}^{T}\right)\left(D^{2} v+V_{x}(\theta+v)\right)\right\rangle \\
& =\left\langle v_{\theta}^{T} D^{2} v\right\rangle+\left\langle M^{T} V_{x}(\theta+v)\right\rangle \tag{50}
\end{align*}
$$

where the latter equality follows since $\left\langle D^{2} v\right\rangle=0$. Let us compute the $i$-th component of $\left\langle v_{\theta}^{T} D^{2} v\right\rangle$. Integrating by parts

$$
\begin{align*}
\left\langle v_{\theta}^{T} D^{2} v\right\rangle_{i} & =\sum_{k=1}^{d} f_{\mathbb{T}^{d}} \frac{\partial v_{k}}{\partial \theta_{i}} D^{2} v_{k} d \theta \\
& =(-1) \sum_{k=1}^{d} f_{\mathbb{T}^{d}} v_{k} \frac{\partial}{\partial \theta_{i}}\left(D^{2} v_{k}\right) d \theta \\
& =(-1)^{2} \sum_{k=1}^{d} f_{\mathbb{T}^{d}}\left(D v_{k}\right) D \frac{\partial v_{k}}{\partial \theta_{i}} d \theta \\
& =(-1)^{3} \sum_{k=1}^{d} f_{\mathbb{T}^{d}}\left(D^{2} v_{k}\right) \frac{\partial v_{k}}{\partial \theta_{i}} d \theta \\
& =-\left\langle v_{\theta}^{T} D^{2} v\right\rangle_{i} . \tag{51}
\end{align*}
$$

Thus $\left\langle v_{\theta}^{T} D^{2} v\right\rangle=0$ and, in view of (50), it remains to check that

$$
\left\langle M^{T} V_{x}(\theta+v)\right\rangle=0
$$

By the chain rule:

$$
\left\langle M^{T} V_{x}(\theta+v)\right\rangle_{i}=\left\langle\sum_{k=1}^{d} \frac{\partial(\theta+v(\theta))_{k}}{\partial \theta_{i}} V_{x_{k}}(\theta+v)\right\rangle=\left\langle\frac{\partial}{\partial \theta_{i}} V(\theta+v)\right\rangle=0
$$

showing that

$$
\begin{equation*}
\left\langle M^{T} \varepsilon\right\rangle=0 \tag{52}
\end{equation*}
$$

Inverting $D$ in (49) we find that

$$
\begin{equation*}
M^{T} M D z=-D^{-1}\left(M^{T} \varepsilon\right)+c \tag{53}
\end{equation*}
$$

where $c$ is a suitable constant vector that we shall shortly identify.
Let

$$
\begin{equation*}
P=M^{T} M \tag{54}
\end{equation*}
$$

and notice that $P=P(\theta)$ is, for $\theta \in \mathbb{T}^{d}$, a strictly positive defined matrix: $P>0$. Thus $P$ is invertible and $P^{-1}>0$. Rephrasing (53) in terms of $P$ :

$$
\begin{equation*}
D z=-P^{-1} D^{-1}\left(M^{T} \varepsilon\right)+P^{-1} c . \tag{55}
\end{equation*}
$$

By the positiveness of $P^{-1}$ and its integrability over $\mathbb{T}^{d}$, we have that also $\left\langle P^{-1}\right\rangle$ is positive and in particular invertible. By taking the average on both sides of (55), we see that in order for (55) to make sense we have to choose:

$$
\begin{equation*}
c:=\left\langle P^{-1}\right\rangle^{-1}\left\langle P^{-1} D^{-1}\left(M^{T} \varepsilon\right)\right\rangle . \tag{56}
\end{equation*}
$$

We can now solve for $z$ obtaining:

$$
\begin{align*}
z & =b+D^{-1}\left(-P^{-1} D^{-1}\left(M^{T} \varepsilon\right)+P^{-1} c\right) \\
& =: b+\hat{z} \tag{57}
\end{align*}
$$

having defined $\hat{z}$ as

$$
\begin{equation*}
\hat{z}=D^{-1}\left(-P^{-1} D^{-1}\left(M^{T} \varepsilon\right)+P^{-1} c\right) \tag{58}
\end{equation*}
$$

and $b$ denotes the arbitrary average of $z$. We fix this ambiguity by requiring that

$$
\begin{equation*}
\left\langle v^{\prime}\right\rangle=\langle v\rangle, \tag{59}
\end{equation*}
$$

which is equivalent to

$$
0=\langle w\rangle=\langle M z\rangle=\langle M b\rangle+\langle M \hat{z}\rangle=\langle M\rangle b+\langle M \hat{z}\rangle=b+\langle M \hat{z}\rangle,
$$

i.e.,

$$
\begin{equation*}
b=-\langle M \hat{z}\rangle . \tag{60}
\end{equation*}
$$

The above analysis may be summarized in the following

Lemma 2 (KAM scheme) Let $V: \mathbb{T}^{d} \rightarrow \mathbb{R}$ be smooth enough and let $\omega \in \mathbb{R}^{d}$ be a diophantine vector. Assume that a smooth enough function $v: \mathbb{T}^{d} \rightarrow \mathbb{R}^{d}$ is given so that

$$
M=I+v_{\theta}
$$

is an invertible matrix on $\mathbb{T}^{d}$ and define

$$
\varepsilon(\theta):=D^{2} v+V_{x}(\theta+v), \quad\left(D:=\sum_{i=1}^{d} \omega_{i} \frac{\partial}{\partial \theta_{i}}\right) .
$$

Then:

$$
\begin{equation*}
\left\langle M^{T} \varepsilon\right\rangle=0 \quad \text { and } \quad\left\langle M^{T} \varepsilon_{\theta}-\varepsilon_{\theta}^{T} M\right\rangle=0 \tag{61}
\end{equation*}
$$

Furthermore, if we let:

$$
\begin{aligned}
P & :=M^{T} M \\
c & :=\langle P\rangle^{-1}\left\langle P^{-1} D^{-1}\left(M^{T} \varepsilon\right)\right\rangle \\
\hat{z} & :=D^{-1}\left(-P^{-1} D^{-1}\left(M^{T} \varepsilon\right)+P^{-1} c\right) \\
b & :=-\langle M \hat{z}\rangle \\
z & :=b+\hat{z} \\
w & :=M z \\
Q_{1} & :=V_{x}(\theta+v+w)-V_{x}(\theta+v)-V_{x x}(\theta+v) w \\
Q_{2} & :=Q_{1}+\varepsilon_{\theta} z \\
Q_{3} & :=Q_{2}+M^{-T}\left(D^{-1}\left(M^{T} \varepsilon_{\theta}-\varepsilon_{\theta}^{T} M\right)\right) D z \\
v^{\prime} & :=v+w \\
\varepsilon^{\prime}(\theta) & :=D^{2} v^{\prime}+V_{x}\left(\theta+v^{\prime}\right)
\end{aligned}
$$

then:

$$
\begin{equation*}
\varepsilon^{\prime}=Q_{3} \quad \text { and } \quad\left\langle v^{\prime}\right\rangle=\langle v\rangle . \tag{62}
\end{equation*}
$$

Remark 7 (i) The above lemma does not contain any quantitative statement, nor its proof uses in any way the fact that $\varepsilon$ should be a "small" function.
(ii) The proof of Lemma 2 is based upon a series of identities: $(50) \div(52)$ and (43) $\div(45)$ [proof of (61)]; (32), (37), (40), (41), (46), (48), (49), (59) [proof of (62)].
(iii) At this level, the above KAM scheme is purely "algebraic" and it will be only after having equipped it with quantitative estimates that it will be possible to iterate the scheme and to actually construct solutions of (23).

### 1.4 Banach spaces of analytic functions and technical lemmata

In this section we introduce "monotone families" of Banach spaces of realanalytic functions on $\mathbb{T}^{d}$; such families will depend upon a parameter $\xi \geq 0$
and "monotone" means that a space parameterized by $\xi>\xi^{\prime}$ is smaller than the space parameterized by $\xi^{\prime}$.

Usually in KAM theory one works either with complex sup-norms or with Fourier norms. In connection with smooth theory (chapter two below) supnorms are more convenient, while for the extension of KAM theory to infinite dimensions Fourier norms are more suited. In this section we shall discuss sup-norms and for completeness we present the analogous technical results also for Fourier norms in Appendix B.

In these lectures we use the following standard notation: for $n \in \mathbb{Z}^{d}, \alpha \in \mathbb{N}^{d}$ and $x \in \mathbb{C}^{d}$, we let

$$
\begin{equation*}
|n|=\sum_{i=1}^{d}\left|n_{i}\right|, \quad|\alpha|=\sum_{i=1}^{d} \alpha_{i}, \quad \partial^{\alpha} f:=\frac{\partial^{\alpha_{1}+\ldots+\alpha_{d}} f}{\partial x_{1}^{\alpha_{1}} \ldots \partial x_{d}^{\alpha_{d}}}, \quad x^{\alpha}:=x_{1}^{\alpha_{1}} \ldots x_{d}^{\alpha_{d}} ; \tag{63}
\end{equation*}
$$

Fourier coefficients of a periodic function will be denoted $f_{n}$; denote, also, by $\Delta_{\xi}^{d}$ the complex strip

$$
\begin{equation*}
\Delta_{\xi}^{d}:=\left\{x \in \mathbb{C}^{d}:\left|\operatorname{Im} x_{j}\right|<\xi, j=1, \ldots, d\right\} \tag{64}
\end{equation*}
$$

For $\xi \geq 0$ we define ${ }^{11}$
$\mathcal{R}_{\xi}\left(\mathbb{T}^{d}, \mathbb{R}^{N}\right):=\left\{f \in C\left(\mathbb{T}^{d}, \mathbb{R}^{N}\right)\right.$ with bounded analytic extension on $\left.\Delta_{\xi}^{d}\right\}$.
$\mathcal{R}_{\xi}\left(\mathbb{T}^{d}, \mathbb{R}^{N}\right)$ endowed with the sup-norm

$$
\|f\|_{\xi}:=\sup _{\Delta_{\xi}^{d}}|f|
$$

is a Banach space.

Remark 8 In the following, we will consider function $f \in \mathcal{R}_{\xi}\left(\mathbb{T}^{d}, X\right)$ with values in a matrix or tensor space $X$; in such cases the definition of the norm will be adapted in the obvious way ${ }^{12}$.

[^7]We proceed to discuss, from a quantitative point of view, the equation $D u=$ $f$ (for $f$ with $\langle f\rangle=0$ ) for $\omega$ diophantine. While it is elementary to get a bound of the form ${ }^{13}$

$$
\|u\|_{\xi-\delta} \leq c(d, \tau) \frac{1}{\gamma \delta^{a}}\|f\|_{\xi}
$$

(with $0<\delta \leq \xi$ ) for some $a>0$, to get the optimal dependence on the "analyticity loss" $\delta$ (i.e., the best $a$ ) is a subtle matter, which was solved by H. Rüssmann. We present a version of Rüssmann's result due to J. Moser (compare also Salamon's paper [2]).

Lemma 3 (Rüssmann, Moser) Let $d \geq 2, \gamma>0$ and $\tau \geq d-1$; let $\omega \in \mathbb{R}^{d}$ be $(\gamma, \tau)$-diophantine and let $f \in \mathcal{R}_{\xi}$ be such that $\langle f\rangle=0$. Denote (as above) by $u:=D^{-1} f$ the unique solution of $D u=f$ with zero average. Then, there exists a constant $c=c(\tau, d)>0$ such that for every

$$
0<\delta \leq \min \{1, \xi\}
$$

one has

$$
\begin{equation*}
\|u\|_{\xi-\delta}=\left\|D^{-1} f\right\|_{\xi-\delta} \leq \frac{c}{\gamma \delta^{\tau}}\|f\|_{L^{2}, \Delta_{\xi}^{d}} \tag{65}
\end{equation*}
$$

where

$$
\|f\|_{L^{2}, \Delta_{\xi}^{d}}:=\sup \left\{\left(\int_{\mathbb{T}^{d}}|f(u+i v)|^{2} d u\right)^{\frac{1}{2}}, \quad|v|<\xi\right\} .
$$

Remark 9 Clearly $\|f\|_{L^{2}, \Delta_{\xi}^{d}} \leq\|f\|_{\xi}$.
In order to prove the above lemma, we shall make use of the following general estimates, the proof of which are deferred to the Appendix A.

Lemma 4 Let $f \in \mathcal{R}_{\xi}$. Then, for every $n \in \mathbb{Z}^{d}$,

$$
\begin{align*}
\left|f_{n}\right| & \leq\|f\|_{L^{2}, \Delta_{\xi}^{d}} e^{-|n| \xi}  \tag{66}\\
& \leq\|f\|_{\xi} e^{-|n| \xi} .
\end{align*}
$$

Furthermore, there exists a constant $c_{0}=c_{0}(d)>0$ such that, for every positive number $\delta<\min \{1, \xi\}$ and for every $x \in \Delta_{\xi-\delta}^{d}$, one has:

$$
\begin{equation*}
|f(x)| \leq \sum_{n \in \mathbb{Z}^{d}}\left|f_{n}\right| e^{-(n \cdot \operatorname{Im} x)} \leq \frac{c_{0}}{\delta^{\frac{d}{2}}}\|f\|_{L^{2}, \Delta_{\xi}^{d}} . \tag{67}
\end{equation*}
$$

[^8]Proof of Lemma 3 As already discussed above, the unique analytic solutions with vanishing mean value of $D u=f$ is given by

$$
u(\theta):=\left(D^{-1} f\right)(\theta)=\sum_{n \in \mathbb{Z}^{d} \backslash\{0\}} u_{n} e^{i(n \cdot \theta)}, \quad u_{n}:=\frac{f_{n}}{i(n \cdot \omega)}
$$

In order to establish the inequality (65) we first single out the subset

$$
J_{0}:=\left\{n \in \mathbb{Z}^{d} \backslash\{0\}:|n \cdot \omega| \geq \frac{\gamma}{2}\right\}
$$

and define

$$
u^{0}(x)=\sum_{n \in J_{0}} u_{n} e^{i(n \cdot x)}
$$

By Lemma 4, we get that, for $|\operatorname{Im} x|<\xi-\delta$ :

$$
\begin{align*}
\left|u^{0}(x)\right| & \leq \sum_{n \in J_{0}}|n \cdot \omega|^{-1}\left|f_{n}\right| e^{-(n \cdot \operatorname{Im} x)} \\
& \leq \frac{2}{\gamma} \sum_{n \in Z^{d}}\left|f_{n}\right| e^{-(n \cdot \operatorname{Im} x)} \\
& \leq \frac{c_{1}}{\gamma \delta^{\tau}}\|f\|_{L^{2}, \Delta_{\xi}^{d}} \tag{68}
\end{align*}
$$

where $c_{1}:=2 c_{0}$ and we have used that $\tau \geq d-1 \geq \frac{d}{2}$.
The more delicate part of the estimate concerns the integer vectors in $\mathbb{Z}^{d} \backslash J_{0}$. First of all, let us assume, without loss of generality, that

$$
\begin{equation*}
\left|\omega_{k}\right|<\left|\omega_{d}\right|, \quad \forall 1 \leq k \leq d-1 \tag{69}
\end{equation*}
$$

let us also introduce the following notation: if $y=\left(y_{1}, \ldots, y_{d}\right)$ is a vector with $d$ components, we denote by $\hat{y}:=\left(y_{1}, \ldots, y_{d-1}\right)$ the vector formed by the first $(d-1)$ components of $y$. Let, now, $K \geq 1$ be a fixed number and for $\nu=1,2 \ldots$, define

$$
\begin{aligned}
J(\nu, K) & :=\left\{n \in \mathbb{Z}^{d}, 0<|n| \leq K:\right. & & \left.2^{\nu} \gamma^{-1}<|n \cdot \omega|^{-1} \leq 2^{\nu+1} \gamma^{-1}\right\} \\
& =\left\{n \in \mathbb{Z}^{d}, 0<|n| \leq K:\right. & & \left.2^{-(\nu+1)} \gamma \leq|n \cdot \omega|<2^{-\nu} \gamma\right\}
\end{aligned}
$$

Here is a list of properties of $\omega$ and $J(\nu, K)$ :
(i) $\left|\omega_{d}\right|>\left|\omega_{k}\right| \geq \gamma$;
(ii) if $n \in J(\nu, K)$ then $\hat{n} \neq 0$;
(iii) if $n, n^{\prime} \in J(\nu, K)$ and $\hat{n}=\hat{n}^{\prime}$ then $n=n^{\prime}$;
(iv) if $n \in \mathbb{Z}^{d}$ is such that $\hat{n} \neq 0$ then

$$
\begin{equation*}
|n \cdot \omega| \geq \frac{\gamma}{(3|\hat{n}|)^{\tau}} \tag{70}
\end{equation*}
$$

(v) if $n, n^{\prime} \in J(\nu, K)$ and $n \neq n^{\prime}$ then $\left|\hat{n}-\hat{n}^{\prime}\right| \geq \frac{2^{\frac{\nu-1}{\tau}}}{3}$;
(vi) there exists a constants $c_{2}=c_{2}(d)>0$ such that

$$
\operatorname{Card} J(\nu, K) \leq c_{2} K^{d-1} 2^{-\frac{\nu(d-1)}{\tau}}
$$

(vii) $J(\nu, K)=\emptyset$ when $2^{\frac{\nu}{\tau}} \geq K$;
(viii) there exists a constant $c_{3}=c_{3}(d)>0$ such that the following holds. If $J(K)$ denotes the set

$$
J(K):=\left\{n \in \mathbb{Z}^{d}: 0<|n| \leq K \text { and }|n \cdot \omega|<\gamma / 2\right\}=\bigcup_{\nu \geq 1} J(\nu, K)
$$

then

$$
\begin{equation*}
\sum_{n \in J(K)} \frac{1}{|n \cdot \omega|} \leq c_{3} \frac{K^{\tau}}{\gamma} \tag{71}
\end{equation*}
$$

## Proof of properties (i) $\div($ (viii)

(i): The first inequality is (69). The second inequality follows from the Diophantine property (27) by taking $n=e_{k}$ (the unit versor in $\mathbb{Z}^{d}$ ).
(ii): If $n \in J(\nu, K)$, then

$$
\begin{equation*}
|n \cdot \omega|<2^{-\nu} \gamma \leq \gamma / 2 ; \tag{72}
\end{equation*}
$$

thus from $\hat{n}=0$ it would follow, by (i), that $|n \cdot \omega|=\left|n_{d} \omega_{d}\right| \geq\left|\omega_{d}\right|>\gamma$, which would contradict (72).
(iii): Assume (by contradiction) that $n, n^{\prime} \in J(\nu, K)$ with $\hat{n}=\hat{n}^{\prime}$ and $n_{d} \neq n_{d}^{\prime}$. Then, by (i) and (72):

$$
\gamma<\left|\omega_{d}\right| \leq\left|n_{d}-n_{d}^{\prime}\right|\left|\omega_{d}\right|=\left|n \cdot \omega-n^{\prime} \cdot \omega\right| \leq|n \cdot \omega|+\left|n^{\prime} \cdot \omega\right|<\frac{\gamma}{2}+\frac{\gamma}{2}=\gamma,
$$

which is impossible.
(iv) Fix $\hat{n} \neq 0$ and choose $n_{d} \in \mathbb{Z}$ so that $n_{\min }:=\left(\hat{n}, n_{d}\right)$ minimizes $|n \cdot \omega|$. Clearly, $n_{\min }$ minimizes also $\left|\hat{n} \cdot \frac{\hat{\omega}}{\omega_{d}}+n_{d}\right|=|n \cdot \omega| /\left|\omega_{d}\right|$. Thus, $\left|\hat{n} \cdot \frac{\hat{\omega}}{\omega_{d}}+n_{d}\right| \leq 1$. Therefore, by (69), $\left|n_{d}\right| \leq 1+\left|\hat{n} \cdot \frac{\hat{\omega}}{\omega_{d}}\right| \leq 1+|\hat{n}| \leq 2|\hat{n}|$, which implies that $\left|n_{\text {min }}\right|=|\hat{n}|+\left|n_{d}\right| \leq 3|\hat{n}|$. In conclusion, by (27), and the above estimates,

$$
|n \cdot \omega| \geq\left|n_{\min } \cdot \omega\right| \geq \frac{\gamma}{\left|n_{\min }\right|^{\tau}} \geq \frac{\gamma}{(3|\hat{n}|)^{\tau}}
$$

(v): By (iii), $\hat{n} \neq \hat{n}^{\prime}$. Thus, by (iv) (applied to the difference $n-n^{\prime}$ ) and by the definition of $J(\nu, K)$, we find $\left|\hat{n}-\hat{n}^{\prime}\right|^{-\tau} \leq \frac{3^{\tau}}{\gamma}\left|\left(n-n^{\prime}\right) \cdot \omega\right| \leq \frac{3^{\tau}}{\gamma}\left(|n \cdot \omega|+\left|n^{\prime} \cdot \omega\right|\right) \leq 3^{\tau} 2^{-(\nu-1)}$, proving the claim.
(vi): By (iii), $J(\nu, K)$ is in a one-to-one correspondence with $\hat{J}(\nu, K):=\left\{\hat{n} \in \mathbb{Z}^{d-1}\right.$ : $n \in J(\nu, K)\}$. By the estimate in (v), the distance between two points in $\hat{J}(\nu, K)$ is at least $\frac{2^{\frac{\nu-1}{\tau}}}{3}$. Thus a simple geometrical argument yields the desired upper bound on the cardinality of $\hat{J}(\nu, K)$ and hence on the cardinality of $J(\nu, K)$.
(vii): If $n \in J(\nu, K)$, by definition $|n \cdot \omega|<\gamma 2^{-\nu}$; on the other hand the Diophantine property (27) implies that $|n \cdot \omega| \geq \gamma|n|^{-\tau} \geq \gamma K^{-\tau}$ implying that $K>2^{\nu / \tau}$, which is equivalent to the claim.
(viii): In view of (vii), $J(K)=\bigcup_{\nu=1}^{\nu_{*}} J(\nu, K)$ where $\nu_{*}$ denotes the integer part of $\tau \log K / \log 2$ (i.e., $\nu_{*}$ is the maximal integer $\nu$ for which $2^{\nu / \tau} \leq K$ : for $\nu>\nu_{*}, J(\nu, K)=$ $\emptyset)$. Thus, by (vii), the definition of $J(\nu, K)$, (vi), one finds

$$
\begin{aligned}
\sum_{n \in J(K)} \frac{1}{|n \cdot \omega|} & \leq \sum_{\nu=1}^{\nu_{*}} \sum_{n \in J(\nu, K)} \frac{1}{|n \cdot \omega|} \\
& \leq \sum_{\nu=1}^{\nu_{*}} \sum_{n \in J(\nu, K)} \frac{2^{\nu+1}}{\gamma}=\sum_{\nu=1}^{\nu_{*}} \frac{2^{\nu+1}}{\gamma} \operatorname{Card} J(\nu, K) \\
& \leq \frac{2 c_{2} K^{d-1}}{\gamma} \sum_{\nu=1}^{\nu_{*}} 2^{\nu \frac{(\tau+1-d)}{\tau}} \\
& \leq \frac{c_{3}}{\gamma} K^{\tau}
\end{aligned}
$$

We are now ready to conclude the proof of (65).

$$
\begin{aligned}
\left\|u-u^{0}\right\|_{\xi-\delta} & \leq \sum_{n \notin J_{0}}\left|f_{n} \| n \cdot \omega\right|^{-1} e^{|n|(\xi-\delta)} \\
& \stackrel{(*)}{\leq}\|f\|_{L^{2}, \Delta_{\xi}^{d}} \sum_{n \notin J_{0}}|n \cdot \omega|^{-1} e^{-|n| \delta} \\
& =\|f\|_{L^{2}, \Delta_{\xi}^{d}} \sum_{k=1}^{\infty} \sum_{n \notin J_{0},|n|^{2}=k}|n \cdot \omega|^{-1} e^{-\sqrt{k} \delta}
\end{aligned}
$$

$$
\begin{aligned}
& =\|f\|_{L^{2}, \Delta_{\xi}^{d}} \sum_{k=1}^{\infty} \sum_{J(\sqrt{k}) \backslash J(\sqrt{k-1})}|n \cdot \omega|^{-1} e^{-\sqrt{k} \delta} \\
& =\|f\|_{L^{2}, \Delta_{\xi}^{d}} \sum_{k=1}^{\infty} \sum_{n \in J(\sqrt{k})}|n \cdot \omega|^{-1}\left(e^{-\sqrt{k} \delta}-e^{-\sqrt{k+1} \delta}\right) \\
& \stackrel{(* *)}{\leq}\|f\|_{L^{2}, \Delta_{\xi}^{d}}\left(\sum_{k=1}^{\infty} \frac{\delta}{2 \sqrt{k}} e^{-\sqrt{k} \delta}\right)\left(\sum_{n \in J(\sqrt{k})}|n \cdot \omega|^{-1}\right) \\
& \leq \frac{c_{3}}{2 \gamma}\|f\|_{L^{2}, \Delta_{\xi}^{d}} \sum_{k=1}^{\infty} k^{\frac{\tau-1}{2}} \delta e^{-\sqrt{k} \delta} \\
& \leq \frac{c_{3}}{2 \gamma} \frac{1}{\delta^{\tau}}\|f\|_{L^{2}, \Delta_{\xi}^{d}} \sup _{0<\lambda \leq 1} \sum_{k=1}^{\infty} \lambda^{\tau+1} k^{\frac{\tau-1}{2}} e^{-\sqrt{k} \lambda} \\
& \pm \frac{c_{4}}{\gamma \delta^{\tau}}\|f\|_{L^{2}, \Delta_{\xi}^{d}},
\end{aligned}
$$

where: $(*)$ is by $(66) ;(* *)$ follows from the elementary bounds $e^{-s}-e^{-s-\varepsilon} \leq$ $\varepsilon e^{-s}$ (any $\left.s>0, \varepsilon>0\right)$ and $\sqrt{t+1}-\sqrt{t} \leq(2 \sqrt{t})^{-1}$ (any $\left.t>0\right) ;(\dagger)$ is by property (viii) above; ( $\ddagger$ ) holds for a suitable constant $c_{4}=c_{4}(d, \tau)>0$ since $\sup _{0<\lambda \leq 1} \sum_{k=1}^{\infty} \lambda^{\tau+1} k^{\frac{\tau-1}{2}} e^{-\sqrt{k} \lambda}<\infty$. The proof is completed if one takes $c=\max \left\{c_{1}, c_{4}\right\}$.

Another fundamental tool are the so-called Cauchy estimates, i.e., the estimates of the sup-norm of derivatives of $f \in \mathcal{R}_{\xi}$ in $\Delta_{\xi-\delta}^{d}$.

Denote by $D_{r}^{d}(x)$ the complex polydisc

$$
D_{r}^{d}(x)=\left\{\theta \in \mathbb{C}^{d}:\left|\theta_{i}-x_{i}\right| \leq r, \forall i\right\}
$$

and let $r>0$ be such that ${ }^{14} D_{r}^{d}(x) \subset \Delta_{\xi}^{d}$. Then, by Cauchy Integral Formula,

$$
\partial^{\alpha} f(x)=\frac{\alpha!}{(2 \pi i)^{d}} \int_{\partial D_{r}^{d}(x)} \frac{f(\theta)}{\left(\theta_{1}-x_{1}\right)^{\alpha_{1}+1} \ldots\left(\theta_{d}-x_{d}\right)^{\alpha_{d}+1}} d \theta .
$$

In particular, taking $r=\delta$, we have for every $x \in \Delta_{\xi-\delta}^{d}$ :

[^9]\[

$$
\begin{aligned}
\left|\partial^{\alpha} f(x)\right| & \leq \frac{\alpha!}{(2 \pi)^{d}} \int_{\partial D_{\delta}^{d}} \frac{|f(\theta)|}{\left|\theta_{1}-x_{1}\right|^{\alpha_{1}+1} \ldots\left|\theta_{d}-x_{d}\right|^{\alpha_{d}+1}} d \theta \\
& \leq \frac{\alpha!}{(2 \pi)^{d}} \frac{\|f\|_{\xi}}{\delta^{|\alpha|+d}} \int_{\partial D_{\delta}^{d}} d \theta \\
& =\frac{\alpha!}{\delta^{|\alpha|}}\|f\|_{\xi} .
\end{aligned}
$$
\]

Thus, the following Cauchy estimate holds

$$
\begin{equation*}
\left\|\partial^{\alpha} f\right\|_{\xi-\delta} \leq \alpha!\delta^{-|\alpha|}\|f\|_{\xi} \tag{73}
\end{equation*}
$$

Combining the above estimates one gets easily the following
Lemma 5 Let $f \in \mathcal{R}_{\xi}\left(\mathbb{T}^{d}, X\right)$, let $p \in \mathbb{N}, \alpha \in \mathbb{N}^{d}$ and assume $\langle f\rangle=0$ when $\alpha=0$. Let $0<\delta \leq \xi$. Then, there exist $C(p, \alpha, \tau)>0$ such that

$$
\left\|D^{-p} \partial^{\alpha} f\right\|_{\xi-\delta} \leq \frac{C(p, \alpha, \tau)}{\gamma^{p} \delta^{p \tau+|\alpha|}}\|f\|_{\xi}
$$

Exercise Give an explicit estimate of $C(p, \alpha, \tau)$.

Remark 10 We shall use also the following trivial facts:
(i) Let $f \in \mathcal{R}_{\xi}\left(\mathbb{T}^{d}, X\right), g \in \mathcal{R}_{\xi}\left(\mathbb{T}^{d}, Y\right)$, with $X$ and $Y$ tensor spaces $\left(\mathbb{R}^{N}\right.$, matrices or higher dimensional tensors) and assume that the product $f g$ is well defined. Then

$$
\|f g\|_{\xi} \leq\|f\|_{\xi}\|g\|_{\xi}
$$

(ii) Let $V \in \mathcal{R}_{\xi_{*}}\left(\mathbb{T}^{d}, \mathbb{R}\right), g \in \mathcal{R}_{\xi}\left(\mathbb{T}^{d}, \mathbb{R}^{d}\right)$, with $\xi_{*}>\xi$. If $\left\|\operatorname{Im} g_{j}\right\|_{\xi} \leq \xi_{*}-\xi$ for all $j=1, \ldots, d$, then $\|V(x+g(x))\|_{\xi} \leq\|V\|_{\xi_{*}}$.

Remark 11 The same estimates could be done also for the non-autonomous case. In this case we consider a function $f(x, t)$ defined on the $(d+1)$ dimensional torus and we require $\omega$ to be $(\gamma, \tau)$-diophantine in the sense that

$$
|\omega \cdot n+m| \geq \frac{\gamma}{|n|^{\tau}}
$$

for some $\tau \geq d, \gamma>0$ and every $(n, m) \in \mathbb{Z}^{d} \times \mathbb{Z}$ with $n \neq 0$.

### 1.5 An analytic KAM theorem

In this section we prove a KAM theorem in the real-analytic setting.

Theorem 6 Fix $0<\bar{\xi}<\xi<\xi_{*} \leq 1$. Let $V \in \mathcal{R}_{\xi_{*}}\left(\mathbb{T}^{d}, \mathbb{R}\right)$, let $v \in$ $\mathcal{R}_{\xi}\left(\mathbb{T}^{d}, \mathbb{R}^{d}\right)$ be such that

$$
\begin{equation*}
\|\operatorname{Im} v\|_{\xi} \leq \xi_{*}-\xi \tag{74}
\end{equation*}
$$

and let $\omega$ be $(\gamma, \tau)$-diophantine ${ }^{15}$. Let, also, $\lambda, \eta$ and $\alpha$ be numbers greater or equal than one such that

$$
\lambda \geq\left\|I+v_{\theta}\right\|_{\xi}, \quad \eta \geq\left\|\left(I+v_{\theta}\right)^{-1}\right\|_{\xi}, \quad \alpha \geq \frac{\left\|V_{x x x}\right\|_{\xi_{*}}}{\gamma^{2}}
$$

There exists a constant $C=C(\tau, d)>1$, such that $i^{16}$

$$
\begin{equation*}
E:=\frac{C}{\gamma^{2}}\|\mathcal{E}(v)\|_{\xi} \alpha(\lambda \eta)^{10}(\xi-\bar{\xi})^{-(4 \tau+2)} \leq 1 \tag{75}
\end{equation*}
$$

then there exists $u \in \mathcal{R}_{\bar{\xi}}\left(\mathbb{T}^{d}, \mathbb{R}^{d}\right)$ with $\langle u\rangle=\langle v\rangle$, which solves the Euler equation

$$
\begin{equation*}
D^{2} u+V_{x}(\theta+u)=: \mathcal{E}(u)=0 \tag{76}
\end{equation*}
$$

Furthermore, there exists a constant $K=K(\tau, d)>0$ such that

$$
\begin{equation*}
\max \left\{\|u-v\|_{\bar{\xi}},\left\|\partial_{\theta} u-\partial_{\theta} v\right\|_{\bar{\xi}}\right\} \leq K E \tag{77}
\end{equation*}
$$

Proof As a first step, we equip the KAM scheme described in Lemma 2 with analytical estimates. The second step will be to iterate the procedure controlling the convergence.

Remark 12 Here and below for simplicity we will use the notation "const" to denote finite (different) constants, which depend only on $\tau$ and $d$.

[^10]Define

$$
\varepsilon(\theta):=D^{2} v+V_{x}(\theta+v) .
$$

and let $\mu>0$ be such that

$$
\|\varepsilon\|_{\xi} \leq \mu
$$

Recalling the definition of $M=I+v_{\theta}$, we have, by hypothesis,

$$
\|M\|_{\xi} \leq \lambda, \quad\left\|M^{-1}\right\|_{\xi} \leq \eta
$$

$\operatorname{Fix}^{17} \bar{\xi}<\xi^{\prime}<\xi$ and let $\delta=\left(\xi-\xi^{\prime}\right) / 2$ :

$$
\begin{equation*}
\xi^{\prime}=: \xi-2 \delta, \quad \delta=\frac{\xi-\xi^{\prime}}{2} \tag{78}
\end{equation*}
$$

and let us denote, as above,

$$
P=M^{T} M
$$

Then

$$
\begin{equation*}
\|P\|_{\xi} \leq \lambda^{2}, \quad\left\|P^{-1}\right\|_{\xi} \leq \eta^{2} \tag{79}
\end{equation*}
$$

We start by estimating $c$ in (56). To estimate it, we will use the fact ${ }^{18}$ that for each positive symmetric matrix $T: \mathbb{T}^{d} \rightarrow \operatorname{Mat}(d \times d)$,

$$
\begin{equation*}
\left\|\langle T\rangle^{-1}\right\| \leq \sup _{\theta \in \mathbb{T}^{d}}\left\|T^{-1}\right\| \tag{80}
\end{equation*}
$$

Therefore, by (80) and (79),

$$
\left\|\left\langle P^{-1}\right\rangle^{-1}\right\| \leq\|P\|_{0} \leq \lambda^{2}
$$

By Lemma 5 we get

$$
\left\|\left\langle P^{-1} D^{-1}\left(M^{T} \varepsilon\right)\right\rangle\right\|_{0} \leq\left\|P^{-1}\right\|_{0}\left\|D^{-1}\left(M^{T} \varepsilon\right)\right\|_{0} \leq \frac{\mathrm{const}}{\gamma \xi^{\tau}} \eta^{2}\left\|M^{T} \varepsilon\right\|_{\xi},
$$

which leads to

$$
\begin{equation*}
|c| \leq \operatorname{const} \frac{\lambda^{3} \eta^{2}}{\gamma \xi^{\tau}} \mu \tag{81}
\end{equation*}
$$

[^11]We proceed to estimate $\hat{z}$ (see (58)). Applying Lemma 5 twice, we have

$$
\begin{aligned}
\|\hat{z}\|_{\xi^{\prime}} & \leq \text { const } \frac{1}{\gamma \delta^{\tau}}\left\|-P^{-1} D^{-1}\left(M^{T} \varepsilon\right)+P^{-1} c\right\|_{\xi-\delta} \\
& \leq \text { const } \frac{\eta^{2}}{\gamma \delta^{\tau}}\left[\left\|D^{-1}\left(M^{T} \xi\right)\right\|_{\xi-\delta}+\frac{\lambda^{3} \eta^{2} \mu}{\gamma \xi^{\tau}}\right] \\
& \leq \text { const } \frac{\eta^{2}}{\gamma \delta^{\tau}}\left[\frac{\lambda^{3} \eta^{2} \mu}{\gamma \xi^{\tau}}+\frac{\lambda \mu}{\gamma \delta^{\tau}}\right] \\
& \leq \text { const } \frac{\lambda^{3} \eta^{4} \mu}{\gamma^{2} \delta^{2 \tau}} .
\end{aligned}
$$

Thus,

$$
\begin{equation*}
\|\hat{z}\|_{\xi^{\prime}} \leq \text { const } \frac{\lambda^{3} \eta^{4}}{\gamma^{2} \delta^{2 \tau}} \mu \tag{82}
\end{equation*}
$$

Now, since $z=b+\hat{z}$ and $b=-\langle M \hat{z}\rangle$, we have

$$
|b|=|\langle M \hat{z}\rangle| \leq \lambda\|\hat{z}\|_{\xi^{\prime}} .
$$

Therefore,

$$
\begin{equation*}
\|z\|_{\xi^{\prime}} \leq \text { const } \frac{(\lambda \eta)^{4}}{\gamma^{2} \delta^{2 \tau}} \mu, \quad\|w\|_{\xi^{\prime}} \leq \text { const } \frac{\lambda^{5} \eta^{4}}{\gamma^{2} \delta^{2 \tau}} \mu \tag{83}
\end{equation*}
$$

Let us estimate now the remainder $Q_{1}$. Using again the standard formula for the remainder of the Taylor expansion and applying Remark 10 and the definition of $\alpha$, we have

$$
\begin{equation*}
\left\|Q_{1}\right\|_{\xi^{\prime}} \leq \frac{1}{2}\left\|V_{x x x}\right\|_{\xi_{*}}\|w\|_{\xi^{\prime}}^{2} \leq \mathrm{const}\left\|V_{x x x}\right\|_{\xi_{*}} \frac{\lambda^{10} \eta^{8} \mu^{2}}{\gamma^{4} \delta^{4 \tau}} \leq \mathrm{const} \frac{\alpha \lambda^{10} \eta^{8}}{\gamma^{2} \delta^{4 \tau}} \mu^{2} \tag{84}
\end{equation*}
$$

provided (compare Remark 10)

$$
\|\operatorname{Im}(v+s w)\|_{\xi^{\prime}} \leq \xi_{*}-\xi^{\prime}
$$

for every $s \in[0,1]$. But, by (74) and (83),

$$
\begin{aligned}
\|\operatorname{Im}(v+s w)\|_{\xi^{\prime}} & \leq\|\operatorname{Im} v\|_{\xi}+\|w\|_{\xi^{\prime}} \leq \xi_{*}-\xi+\|w\|_{\xi^{\prime}} \\
& \leq \xi_{*}-\xi+\mathrm{const} \frac{\lambda^{5} \eta^{4}}{\gamma^{2} \delta^{2 \tau}} \mu \\
& \leq \xi_{*}-\xi^{\prime}
\end{aligned}
$$

which is implied if we assume that

$$
\begin{equation*}
\text { const } \frac{(\lambda \eta)^{5}}{\gamma^{2} \delta^{2 \tau+1}} \mu \leq 1 \tag{85}
\end{equation*}
$$

Since $Q_{2}=Q_{1}+\varepsilon_{\theta} z$, by Lemma 5 and Remark 10, we find

$$
\begin{equation*}
\left\|\varepsilon_{\theta}\right\|_{\xi-\delta} \leq \frac{\mu}{\delta} \tag{86}
\end{equation*}
$$

and

$$
\left\|Q_{2}\right\|_{\xi^{\prime}} \leq\left\|Q_{1}\right\|_{\xi^{\prime}}+\left\|\varepsilon_{\theta} z\right\|_{\xi^{\prime}} \leq\left\|Q_{1}\right\|_{\xi^{\prime}}+\left\|\varepsilon_{\theta}\right\|_{\xi^{\prime}}\|z\|_{\xi^{\prime}} \leq\left\|Q_{1}\right\|_{\xi^{\prime}}+\frac{\mu}{2 \delta}\|z\|_{\xi^{\prime}}
$$

Thus, by (83) and (84),

$$
\begin{equation*}
\left\|Q_{2}\right\|_{\xi^{\prime}} \leq \mathrm{const} \frac{\alpha \lambda^{10} \eta^{8}}{\gamma^{2} \delta^{4 \tau}} \mu^{2} \tag{87}
\end{equation*}
$$

Here we implicitly used the fact that $\tau \geq 1 / 2$ and $\delta \leq 1$.
We turn to the estimate of the norm of $Q_{3}$ (defined in (48)): we will use again Lemma 5. First, observe that $D z=D \hat{z}$ and that

$$
D \hat{z}=-P^{-1} D^{-1}\left(M^{T} \varepsilon\right)+P^{-1} c .
$$

Thus

$$
\|D \hat{z}\|_{\xi^{\prime}}=\left\|-P^{-1} D^{-1}\left(M^{T} \varepsilon\right)+P^{-1} c\right\|_{\xi^{\prime}} \leq \text { const } \frac{\lambda^{3} \eta^{4}}{\gamma \delta^{\tau}} \mu
$$

Recalling (86), we find

$$
\left\|D^{-1}\left(M^{T} \varepsilon_{\theta}-\varepsilon_{\theta}^{T} M\right)\right\|_{\xi^{\prime}} \leq \text { const } \frac{1}{\gamma \delta^{\tau}}\left\|M^{T} \varepsilon_{\theta}\right\|_{\xi-\delta} \leq \mathrm{const} \frac{\lambda}{\gamma \delta^{\tau+1}} \mu
$$

As for the second term in the formula for $Q_{3}$, we find

$$
\left\|M^{-T}\left(D^{-1}\left(M^{T} \varepsilon_{\theta}-\varepsilon_{\theta}^{T} M\right)\right) D z\right\|_{\xi^{\prime}} \leq \mathrm{const} \frac{\lambda \eta \mu}{\gamma \delta^{\tau+1}} \frac{\lambda^{3} \eta^{4} \mu}{\gamma \delta^{\tau}} \leq \mathrm{const} \frac{\lambda^{4} \eta^{5}}{\gamma^{2} \delta^{2 \tau+1}} \mu^{2}
$$

Finally, recalling (87), we get the following bound

$$
\begin{equation*}
\left\|Q_{3}\right\|_{\xi^{\prime}} \leq \text { const } \frac{\alpha \lambda^{10} \eta^{8}}{\gamma^{2} \delta^{4 \tau}} \mu^{2} \leq \mathrm{const} \frac{\alpha(\lambda \eta)^{10}}{\gamma^{2} \delta^{4 \tau}} \mu^{2} \tag{88}
\end{equation*}
$$

Our next step will be to estimate $v^{\prime}, M^{\prime}:=M+w_{\theta}$ and $\varepsilon^{\prime}$ in terms of $v, M$ and $\varepsilon$.

Let us start with $w_{\theta}$. Here again, as we have already done for $\hat{z}$, we will apply Lemma 5 twice. In fact,

$$
\left\|w_{\theta}\right\|_{\xi^{\prime}} \leq \frac{1}{\delta}\|w\|_{\xi-\delta}
$$

Then, since $\xi^{\prime}=\xi-\delta=\xi-\delta / 2-\delta / 2$, using twice Lemma 5, we get:

$$
\begin{aligned}
\|\hat{z}\|_{\xi-\delta} & \leq \operatorname{const} \frac{1}{\gamma \delta^{\tau}}\left\|-P^{-1} D^{-1}\left(M^{T} \varepsilon\right)+P^{-1} c\right\|_{\xi-\delta / 2} \\
& \leq \operatorname{const} \frac{\eta^{2}}{\gamma \delta^{\tau}}\left[\frac{\lambda^{3} \eta^{2} \mu}{\gamma \xi^{\tau}}+\left\|D^{-1}\left(M^{T} \varepsilon\right)\right\|_{\xi-\delta / 2}\right] \\
& \leq \operatorname{const} \frac{\eta^{2}}{\gamma \delta^{\tau}}\left[\frac{\lambda^{3} \eta^{2} \mu}{\gamma \xi^{\tau}}+\frac{\lambda \mu}{\gamma(\delta / 2)^{\tau}}\right] \\
& \leq \operatorname{const} \frac{\eta^{4} \lambda^{3}}{\gamma^{2} \delta^{2 \tau}} \mu
\end{aligned}
$$

Since $z=b+\hat{z}$ and $w=M z$, we get the estimate

$$
\begin{equation*}
\left\|w_{\theta}\right\|_{\xi^{\prime}} \leq \mathrm{const} \frac{\lambda^{5} \eta^{4}}{\gamma^{2} \delta^{2 \tau+1}} \mu \tag{89}
\end{equation*}
$$

Next, using (89), we easily get

$$
\begin{align*}
\left\|M^{\prime}\right\|_{\xi^{\prime}} & \leq\left\|M+w_{\theta}\right\|_{\xi^{\prime}} \leq \lambda+\text { const } \frac{\lambda^{5} \eta^{4} \mu}{\gamma^{2} \delta^{2 \tau+1}} \\
& =\lambda\left(1+\mathrm{const} \frac{(\lambda \eta)^{4} \mu}{\gamma^{2} \delta^{2 \tau+1}}\right)=: \lambda^{\prime} \tag{90}
\end{align*}
$$

As for the inverse matrix

$$
\begin{aligned}
\left\|\left(M^{\prime}\right)^{-1}\right\|_{\xi^{\prime}} & \leq\left\|\left(M+w_{\theta}\right)^{-1}\right\|_{\xi^{\prime}}=\left\|\left(I+M^{-1} w_{\theta}\right)^{-1} M^{-1}\right\|_{\xi^{\prime}} \\
& \leq \eta\left(1-\mathrm{const} \frac{(\lambda \eta)^{5} \mu}{\gamma^{2} \delta^{2 \tau+1}}\right)^{-1},
\end{aligned}
$$

provided

$$
\text { const } \frac{(\lambda \eta)^{5} \mu}{\gamma^{2} \delta^{2 \tau+1}}<1
$$

In fact, assuming that

$$
\begin{equation*}
\text { const } \frac{(\lambda \eta)^{5}}{\gamma^{2} \delta^{2 \tau+1}} \mu \leq \frac{1}{2} \tag{91}
\end{equation*}
$$

one finds ${ }^{19}$

$$
\begin{equation*}
\left\|\left(M^{\prime}\right)^{-1}\right\|_{\xi^{\prime}} \leq \eta\left(1+\mathrm{const} \frac{(\lambda \eta)^{5}}{\gamma^{2} \delta^{2 \tau+1}} \mu\right)=: \eta^{\prime} \tag{92}
\end{equation*}
$$

Note that (91) is the same condition (up to the constant) as (85). Since, by definition, $\varepsilon^{\prime}=Q_{3}$, we have obtained

$$
\begin{equation*}
\left\|\varepsilon^{\prime}\right\|_{\xi^{\prime}} \leq \mathrm{const} \frac{\alpha(\lambda \eta)^{10}}{\gamma^{2} \delta^{4 \tau}} \mu^{2}=: \mu^{\prime} \tag{93}
\end{equation*}
$$

It will be useful to introduce also a "dimensionless" parameter

$$
\begin{equation*}
\bar{\mu}=\mu \gamma^{-2} \tag{94}
\end{equation*}
$$

in terms of which (93) may be rewritten as

$$
\begin{equation*}
\left\|\varepsilon^{\prime}\right\|_{\xi^{\prime}} \gamma^{-2} \leq\left(\text { const } \frac{\sqrt{\alpha}(\lambda \eta)^{5} \bar{\mu}}{\delta^{2 \tau}}\right)^{2} \leq\left(\operatorname{const} \frac{\sqrt{\alpha}(\lambda \eta)^{5} \bar{\mu}}{\delta^{2 \tau+1}}\right)^{2}=: \bar{\mu}^{\prime} \tag{95}
\end{equation*}
$$

From now on we will replace (91) by the stronger condition

$$
\begin{equation*}
\text { const } \frac{\sqrt{\alpha}(\lambda \eta)^{5} \bar{\mu}}{\delta^{2 \tau+1}} \leq 1 \tag{96}
\end{equation*}
$$

Let us now turn to the second step, i.e., to the control of the convergence of the iteration process. For $i \geq 0$, let the input data $v, \varepsilon, \xi$ correspond to the $i$-th step of iteration, and let the output $v^{\prime}, \varepsilon^{\prime}$ and $\xi^{\prime}$ correspond to $(i+1)$-th step: in particular the function $v$ and the parameter $\xi$ in the statement of Theorem 6 will be denoted, respectively, $v_{0}, \xi_{0}$. Thus,

$$
v_{i}=v_{0}+\sum_{j=0}^{i-1} w_{j}
$$

[^12]and our aim is to show that this sequence converges to some real-analytic function $u$, which solves the Euler equation
$$
D^{2} u+V_{x}(\theta+u(t))=0 .
$$

We fix a sequence $\left\{\xi_{i}\right\}$ as follows

$$
\xi_{i}=\bar{\xi}+\frac{\xi-\bar{\xi}}{2^{i}} .
$$

So, $\left\{\xi_{i}\right\}$ is a decreasing sequence, which tends to $\bar{\xi}$. In view of the above definition of $\delta$, we fix also

$$
\delta_{i}:=\frac{\xi_{i}-\xi_{i+1}}{2}=\frac{\xi-\bar{\xi}}{2^{i+2}} .
$$

If (96) holds at each each step of the iteration, i.e., if

$$
\begin{equation*}
\text { const } \frac{\sqrt{\alpha}\left(\lambda_{j} \eta_{j}\right)^{5} \bar{\mu}_{j}}{\delta_{j}^{2 \tau+1}} \leq 1, \quad j=0,1, \ldots, i \tag{97}
\end{equation*}
$$

then, in particular, we see that (compare (90) and (92) attaching the indices $i$ and $i+1$ in the obvious way)

$$
\begin{equation*}
\lambda_{j} \leq 2^{j} \lambda_{0}, \quad \eta_{j} \leq 2^{j} \eta_{0}, \quad j=0,1, \ldots, i \tag{98}
\end{equation*}
$$

In terms of $\lambda_{0}$ and $\eta_{0}$ condition (97) can be rewritten in the form

$$
C_{0} \frac{\sqrt{\alpha}\left(\lambda_{0} \eta_{0}\right)^{5}\left(2^{2 \tau+11}\right)^{i} \bar{\mu}_{i}}{(\xi-\bar{\xi})^{2 \tau+1}} \leq 1
$$

where $C_{0}=C_{0}(d, \tau)$ denotes the largest constant " const " occurred until now. Denoting,

$$
A:=\left(C_{0} \frac{\sqrt{\alpha}\left(\lambda_{0} \eta_{0}\right)^{5}}{(\xi-\bar{\xi})^{2 \tau+1}}\right)^{2}, \quad B:=2^{2 \tau+11}
$$

we see that (95) yields

$$
\begin{equation*}
\bar{\mu}_{i+1} \leq A B^{i} \bar{\mu}_{i}^{2} \leq 1 \tag{99}
\end{equation*}
$$

Such relation may be rewritten as

$$
\hat{\mu}_{i+1} \leq \hat{\mu}_{i}^{2}, \quad \hat{\mu}_{i}:=A B^{i+1} \bar{\mu}_{i}
$$

which, iterated, leads to

$$
\begin{equation*}
\bar{\mu}_{i} \leq \frac{\left(A B \bar{\mu}_{0}\right)^{2^{i}}}{A B^{i+1}} \tag{100}
\end{equation*}
$$

In particular, one can conclude that the iteration process converges if

$$
A B \bar{\mu}_{0}<1, \quad \bar{\mu}_{0} \geq \frac{\left\|\varepsilon_{0}\right\|_{\xi_{0}}}{\gamma^{2}}
$$

showing, in particular, that $v+\sum_{j=0}^{\infty} w_{j}$ converges uniformly on the complex strip of width $\bar{\xi}$ to the real-analytic function

$$
u:=v+\sum_{j=0}^{\infty} w_{j}
$$

which (since $\varepsilon_{j} \rightarrow 0$ uniformly) will satisfy the Euler equation (76).
Finally, we prove (77). First of all note that in fact $E$ is nothing else but

$$
E=A B \bar{\mu}_{0} .
$$

Now, (compare (98)),

$$
\begin{aligned}
\left\|w_{j}\right\|_{\bar{\xi}} & \leq \operatorname{const} \frac{2^{2 \tau(j+2)} \lambda_{j}^{5} \eta_{j}^{4} \bar{\mu}_{j}}{(\xi-\bar{\xi})^{2 \tau}}=\operatorname{const}\left(4^{\tau j} \lambda_{j}^{5} \eta_{j}^{4} \bar{\mu}_{j}\right) \\
& \leq\left(\operatorname{const}\left(\lambda_{0} \eta_{0}\right)^{5}\right) 2^{10 j+2 \tau j} \bar{\mu}^{j}=: \tilde{C} 2^{10 j+2 \tau j} \bar{\mu}_{j}
\end{aligned}
$$

where $\tilde{C}:=$ const $\left(\lambda_{0} \eta_{0}\right)^{5}$. Since by (100)

$$
\bar{\mu}_{i}<\left(A B \mu_{0}\right)^{2^{i}}
$$

we have

$$
\sum_{j=0}^{\infty} 2^{10 j+2 \tau j} \bar{\mu}_{j} \leq \sum_{j=0}^{\infty} 2^{2(5+\tau) j}\left(A B \bar{\mu}_{0}\right)^{2^{j}}=\sum_{j=0}^{\infty}\left(C_{1}^{2}\right)^{j} E^{2^{j}}
$$

where we have denoted $C_{1}=2^{5+\tau}$. Continuing the last inequality, we get

$$
\begin{aligned}
\sum_{j=0}^{\infty} C_{1}^{2 j} E^{2^{j}} & <\sum_{j=0}^{\infty} C_{1}^{2^{j}} E^{2^{j}}=\sum_{j=0}^{\infty}\left(C_{1} E\right)^{2^{j}} \\
& \leq \sum_{j=1}^{\infty}\left(C_{1} E\right)^{j}=\frac{C_{1} E}{1-C_{1} E} \\
& \leq C_{1} E\left(1+2 C_{1} E\right) \leq 2 C_{1} E,
\end{aligned}
$$

provided $0<C_{1} E \leq 1 / 2$. This last assumption can be always satisfied by the right choice of the constant.

We conclude this section with an immediate application of Theorem 6 to the "nearly-integrable" case.

Corollary 1 If

$$
\begin{equation*}
\left\|V_{x}\right\|_{\xi} \leq \frac{\gamma^{2}}{C} \frac{(\xi-\bar{\xi})^{4 \tau+2}}{\max \left\{1,\left\|V_{x x x}\right\|_{\xi_{*}} \gamma^{-2}\right\}} \tag{101}
\end{equation*}
$$

then there exists a function $u \in \mathcal{R}_{\xi}\left(\mathbb{T}^{d}, \mathbb{R}^{d}\right)$ such that $\langle u\rangle=0$ which solves the Euler equation $D^{2} u+V_{x}(\theta+u)=0$ with

$$
\|u\|_{\bar{\xi}} \leq \frac{C}{\gamma^{2}}\left\|V_{x}\right\|_{\xi} \frac{\max \left\{1,\left\|V_{x x x}\right\|_{\xi_{*}} \gamma^{-2}\right\}}{(\xi-\bar{\xi})^{4 \tau+2}}
$$

Proof Take as initial approximate solution the function $v \equiv 0$. Then $\mathcal{E}(v)=$ $\mathcal{E}(0)=V_{x}(\theta)$ and one can take $\lambda=\eta=1$ so that (101) is recognized to be (75).

Remark 13 Let $\left\|V_{x}\right\|_{\xi},\left\|V_{x x x}\right\|_{\xi} \leq \varepsilon$. From the properties of diophantine numbers it follows that, if we denote

$$
\Omega_{r}=\left\{\omega \in B_{r}^{d}: \quad|\omega \cdot n| \geq \frac{\gamma}{|n|^{\tau}} \quad \forall n \in \mathbb{Z}^{d} \backslash\{0\}\right\}
$$

then

$$
\operatorname{meas}\left(B_{r}^{d} \backslash \Omega_{r}\right) \leq \text { const meas }\left(B_{r}^{d}\right) \gamma
$$

Now, condition (101) can be met by taking $\gamma=\sqrt{\varepsilon} \hat{C}$ with $\hat{C}$ big enough, showing that the set of $\omega$ 's for which we can find simultaneously a solution for the Euler equation fills (as $\varepsilon \rightarrow 0$ ) a ball of radius $r$ up to a set of measure at most const $\sqrt{\varepsilon}$.

### 1.6 Local uniqueness

In this section we formulate a sufficient condition which provides "local" uniqueness for the solution of Euler equation. First we remark that if $u$ verifies

$$
\begin{equation*}
D^{2} u+V_{x}(\theta+u)=0 \tag{102}
\end{equation*}
$$

then also

$$
\bar{u}:(\theta) \mapsto c+u(\theta+c)
$$

is a solution of the same equation, for every constant $c \in \mathbb{R}^{d}$. Since $\langle\bar{u}\rangle=$ $\langle u\rangle+c$, it is natural to investigate local uniqueness of solutions with prescribed average.

Proposition 7 Let $\omega \in \mathbb{R}^{d}$ be $(\gamma, \tau)$-diophantine. Let $V \in \mathcal{R}_{\xi_{*}}\left(\mathbb{T}^{d}, \mathbb{R}\right)$ and $u, \bar{u} \in \mathcal{R}_{\xi}\left(\mathbb{T}^{d}, \mathbb{R}^{d}\right)$. Assume that $u$ and $\bar{u}$ are two solutions of (102) such that $\langle u\rangle=\langle\bar{u}\rangle$. Assume moreover that $I+u_{\theta}$ is invertible everywhere on $\mathbb{T}^{d}$ and that

$$
\begin{aligned}
\|u\|_{\xi},\|\bar{u}\|_{\xi} & \leq \xi_{*}-\xi \\
\left\|\left(I+u_{\theta}\right)^{-1}\right\|_{\xi} & \leq \eta<+\infty \\
\left\|I+u_{\theta}\right\|_{\xi} & \leq \lambda<+\infty
\end{aligned}
$$

Define

$$
\begin{equation*}
c:=c_{0} \frac{\gamma^{2} \xi^{2 \tau}}{\lambda^{5} \eta^{4}\left\|V_{x x x}\right\|_{\xi_{*}}} \tag{103}
\end{equation*}
$$

where $c_{0}=c_{0}(d, \tau) \geq 1$ is a suitable constant. Then, if $\|u-\bar{u}\|_{\xi}<c$ one has that $u \equiv \bar{u}$.

Proof Let

$$
w:=\bar{u}-u
$$

and notice that $\langle w\rangle=0$. Since $\bar{u}$ and $u$ are solutions of (102), we have

$$
\begin{align*}
0 & =D^{2} \bar{u}+V_{x}(\theta+\bar{u}) \\
& =D^{2} w+D^{2} u+V_{x}(\theta+w+u) \\
& =D^{2} w+V_{x}(\theta+w+u)-V_{x}(\theta+u) \\
& =D^{2} w+V_{x x}(\theta+u) w+Q \tag{104}
\end{align*}
$$

where

$$
Q=\int_{0}^{1}(1-s) V_{x x x}(\theta+u+s w) w w d s
$$

From the expression of $Q$ and Lemma 15 [Remark 10] we easily get

$$
\begin{equation*}
\|Q\|_{\xi^{\prime}} \leq \frac{1}{2}\left\|V_{x x x}\right\|_{\xi_{*}}\|w\|_{\xi^{\prime}}^{2} \tag{105}
\end{equation*}
$$

for every $\xi^{\prime} \in[0, \xi]$.
Let $M=I+u_{\theta}$. Differentiating with respect to $\theta$ equation (102) we get the equality

$$
\begin{equation*}
V_{x x}(\theta+u)=-\left(D^{2} M\right) M^{-1} \tag{106}
\end{equation*}
$$

that we can plug in (104) obtaining

$$
0=D^{2} w-\left(D^{2} M\right) M^{-1} w+Q
$$

Letting

$$
z=M^{-1} w
$$

we have

$$
\begin{aligned}
0 & =D^{2}(M z)-\left(D^{2} M\right) z+Q \\
& =D(M D z)+D(D M z)-\left(D^{2} M\right) z+Q \\
& =D(M D z)+(D M)(D z)+Q
\end{aligned}
$$

which can be rewritten as

$$
\begin{equation*}
0=M^{-T}\left(M^{T} D(M D z)+M^{T}(D M)(D z)\right)+Q \tag{107}
\end{equation*}
$$

Moreover from (106), we have that the matrix $\left(D^{2} M\right) M^{-1}$ is symmetric and so

$$
\begin{aligned}
0 & =M^{T}\left(D^{2} M\right)-\left(D^{2} M\right)^{T} M \\
& =D\left(M^{T} D M-\left(D M^{T}\right) M\right)
\end{aligned}
$$

in particular, since $D^{-1} 0=0$ we have

$$
M^{T} D M-\left(D M^{T}\right) M=\left\langle M^{T} D M-\left(D M^{T}\right) M\right\rangle
$$

By equation (44) we already know, however, that

$$
\left\langle M^{T} D M-\left(D M^{T}\right) M\right\rangle=0
$$

for every matrix $M$ of the form $I+u_{\theta}$. Thus

$$
M^{T} D M-D M^{T} M=0
$$

i.e.

$$
M^{T} D M=D M^{T} M
$$

From (107) it follows:

$$
\begin{aligned}
0 & =M^{-T}\left(M^{T} D(M D z)+\left(D M^{T}\right) M D z\right)+Q \\
& =M^{-T} D\left(M^{T} M D z\right)+Q
\end{aligned}
$$

which means, setting $P=M^{T} M$, that $P D z=-D^{-1}\left(M^{T} Q\right)+c_{1}$ for a suitable constant vector $c_{1}$. Thus

$$
\begin{equation*}
D z=-P^{-1} D^{-1}\left(M^{T} Q\right)+P^{-1} c_{1} \tag{108}
\end{equation*}
$$

and

$$
\begin{equation*}
w=M D^{-1}\left(-P^{-1} D^{-1}\left(M^{T} Q\right)+P^{-1} c_{1}\right)+M c_{2} \tag{109}
\end{equation*}
$$

for a suitable constant vector $c_{2}$. Taking averages in (108) and (109) we obtain the following expressions for $c_{1}$ and $c_{2}$ :

$$
\begin{align*}
& c_{1}=\left\langle P^{-1}\right\rangle^{-1}\left\langle P^{-1} D^{-1}\left(M^{T} Q\right)\right\rangle  \tag{110}\\
& c_{2}=-\left\langle M D^{-1}\left(-P^{-1} D^{-1}\left(M^{T} Q\right)+P^{-1} c_{1}\right)\right\rangle . \tag{111}
\end{align*}
$$

Let us now define, for every $j \in \mathbb{N}, \xi_{j}=2^{-j} \xi$. For every $j \in \mathbb{N}$, by estimates similar to the ones already seen in the previous section we get from (110) and (111) the following inequalities:

$$
\begin{aligned}
& \left|c_{1}\right| \leq \text { const } \lambda^{3} \eta^{2} \frac{2^{j \tau}}{\gamma \xi^{\tau}}\|Q\|_{\xi_{j}} \\
& \left|c_{2}\right| \leq \text { const } \lambda^{4} \eta^{4} \frac{4^{(j+1) \tau}}{\gamma^{2} \xi^{2 \tau}}\|Q\|_{\xi_{j}}
\end{aligned}
$$

which can be inserted in (109) obtaining

$$
\begin{equation*}
\text { const }\|w\|_{\xi_{j+1}} \leq \lambda^{5} \eta^{4} \frac{4^{2 j \tau}}{\gamma^{2} \xi^{2 \tau}}\|Q\|_{\xi_{j}} \tag{112}
\end{equation*}
$$

Letting

$$
k=\max \left\{1, \text { const } \lambda^{5} \eta^{4} \frac{1}{\gamma^{2} \xi^{2 \tau}}\left\|V_{x x x}\right\|_{\xi_{*}}\right\}
$$

we obtain from (112) and (105)

$$
\|w\|_{\xi_{j+1}} \leq k 4^{2 j \tau}\|w\|_{\xi_{j}}
$$

and, iterating as done above (compare (99), (100)), we get

$$
\|w\|_{0} \leq\|w\|_{\xi_{j+1}} \leq\left(k 4^{2 \tau}\|w\|_{\xi}\right)^{2^{j}}
$$

showing that $\|w\|_{0}=0$ (and hence by analyticity $w \equiv 0$ ) whenever

$$
k 4^{2 \tau}\|w\|_{\xi}<1
$$

## 2 Smooth KAM Theory

The aim of this chapter is to exhibit a result of existence of quasi-periodic solutions for systems that are no more required to be analytic but just smooth enough. We will heavily use the previous results, passing through analytic approximations of smooth functions.

### 2.1 Approximation Theory

Here we prove the necessary technical approximation results.
We start by introducing Hölder norms. First of all, for every $l_{0} \in \mathbb{N}$ and for every $f \in C^{l_{0}}\left(\mathbb{R}^{m}\right)$, we define

$$
|f|_{C^{l_{0}}}=\sup _{|\alpha| \leq l_{0} \mathbb{R}^{m}} \sup \left|\partial^{\alpha} f\right| .
$$

If $l=l_{0}+\mu$ with $l_{0} \in \mathbb{N}$ and $\mu \in(0,1)$, we set

$$
|f|_{C^{l}}:=|f|_{C^{l_{0}}}+\sup _{|\alpha|=l_{0}} \sup _{0<|x-y|<1} \frac{\left|\partial^{\alpha} f(x)-\partial^{\alpha} f(y)\right|}{|x-y|^{\mu}} .
$$

For every $l \geq 0$ we define

$$
C^{l}\left(\mathbb{R}^{d}, \mathbb{R}^{m}\right)=\left\{f: \mathbb{R}^{d} \rightarrow \mathbb{R}^{m}: \quad|f|_{C^{l}}<+\infty\right\}
$$

The space $C^{l}\left(\mathbb{R}^{d}, \mathbb{R}^{m}\right)$ endowed with the norm $|\cdot|_{C^{l}}$ is a Banach space; the subspace of $C^{l}\left(\mathbb{R}^{d}, \mathbb{R}^{m}\right)$ made of functions which are $2 \pi$-periodic in each variable will be denoted $C^{l}\left(\mathbb{T}^{d}, \mathbb{R}^{m}\right)$.

Remark 14 In this section it is convenient to work with Euclidean norms on vectors and the associated operator norms on matrices and tensors.

Proposition 8 (Jackson, Moser, Zehnder) Let $l \geq 0, d \in \mathbb{Z}_{+}$and $f \in$ $C^{l}\left(\mathbb{R}^{d}\right)$. There exists a constant $c=c(l, d)>0$ such that for every $0<r \leq 1$ there exists a real analytic function $f_{r}$ on $\Delta_{r}^{d}$ which satisfies ${ }^{20}$

$$
\begin{equation*}
\left|\partial^{\alpha} f_{r}(x)-\sum_{|\beta| \leq l-|\alpha|} \partial^{\alpha+\beta} f(\operatorname{Re} x) \frac{(i \operatorname{Im} x)^{\beta}}{\beta!}\right| \leq c|f|_{C^{l}} r^{l-|\alpha|}, \quad \forall x \in \Delta_{r}^{d} \tag{113}
\end{equation*}
$$

for all $\alpha$ such that $|\alpha| \leq l$.

[^13]In fact, the analytic extension $f_{r}$ may be defined as follows. Let $\phi_{1}$ be an even function in $C_{0}^{\infty}(\mathbb{R})$ with support $[-1,1]$, increasing in $[-1,0]$ and such that

$$
\phi_{1}(0)=1, \quad \quad \partial^{n} \phi_{1}(0)=0 \quad(\forall n \geq 1) ;
$$

for $\xi \in \mathbb{R}^{d}$, let $\phi(\xi)=\phi_{1}\left(|\xi|^{2}\right)$ and let $K$ be the anti-Fourier transform of $\phi$ :

$$
K(x)=\frac{1}{(2 \pi)^{d}} \int_{\mathbb{R}^{d}} \phi(\xi) e^{i x \cdot \xi} d \xi ;
$$

then $f_{r}$ can be taken to be

$$
\begin{align*}
f_{r}(x) & :=\left(\mathcal{S}_{r} f\right)(x):=r^{-d} \int_{\mathbb{R}^{d}} K\left(\frac{x-y}{r}\right) f(y) d y  \tag{114}\\
& =\int_{\mathbb{R}^{d}} K\left(\frac{x}{r}-\xi\right) f(r \xi) d \xi=\int_{\mathbb{R}^{d}} K(\eta) f(x-r \eta) d \eta .
\end{align*}
$$

Proof Since $\phi$ is a real smooth function with compact support, $K$ is real analytic on $\mathbb{C}^{d}$. Some properties of $K$ are collected in the following

Lemma 9 The derivatives of $K$ satisfy

$$
\begin{gather*}
\forall p \in \mathbb{N}, \quad \exists c_{p}:\left|\partial^{\beta} K(x)\right| \leq c_{p} \frac{e^{|\operatorname{Im} x|}}{(1+|x|)^{p}}, \quad \forall|\beta| \leq p  \tag{115}\\
\sup _{x \in \mathbb{R}^{d}} \sup _{\beta \in \mathbb{N}^{d}}\left|\partial^{\beta} K(x)\right| \leq \frac{1}{(2 \pi)^{d}}\|\phi\|_{L^{1}} . \tag{116}
\end{gather*}
$$

Furthermore, if $\alpha, \beta \in \mathbb{N}^{d}$ and $x=u+i v \in \mathbb{C}^{d}$, then ${ }^{21}$

$$
I_{\alpha, \beta}:=\int_{\mathbb{R}^{d}} u^{\beta} \partial^{\alpha} K(u+i v) d u=\left\{\begin{array}{cc}
(-1)^{|\beta|} \frac{\beta!}{(\beta-\alpha)!}(i v)^{\beta-\alpha}, & \text { if } \alpha \leq \beta  \tag{117}\\
0, & \text { otherwise }
\end{array}\right.
$$

[^14]Proof First of all, remark that if $u \in \operatorname{supp} \phi=B_{1}(0)$, then

$$
\left|e^{i x \cdot u}\right|=e^{-\operatorname{Im} x \cdot u} \leq e^{|\operatorname{Im} x|}
$$

Let us denote $\phi_{\beta}(u)=u^{\beta} \phi(u)$. We have

$$
\begin{aligned}
\partial^{\beta} K(x) & :=\partial^{\beta}\left(\frac{1}{(2 \pi)^{d}} \int_{\mathbb{R}^{d}} \phi(u) e^{i x \cdot u} d u\right) \\
& =\frac{i^{|\beta|}}{(2 \pi)^{d}} \int_{\mathbb{R}^{d}} \phi_{\beta}(u) e^{i x \cdot u} d u
\end{aligned}
$$

and, for any multi-index $\alpha \in \mathbb{N}^{d},|\alpha|$ integrations by part give

$$
x^{\alpha} \partial^{\beta} K(x)=\frac{i^{|\alpha|+|\beta|}}{(2 \pi)^{d}} \int_{\mathbb{R}^{d}} \partial^{\alpha} \phi_{\beta}(u) e^{i x \cdot u} d u
$$

Hence

$$
\left|x^{\alpha}\right|\left|\partial^{\beta} K(x)\right| \leq \frac{1}{(2 \pi)^{d}}\left|\partial^{\alpha} \phi_{\beta}\right|_{L^{1}\left(\mathbb{R}^{d}\right)} e^{|\operatorname{Im} x|}
$$

Now remark that for any $p \in \mathbb{N}$,

$$
(1+|x|)^{p} \leq\left(1+\left|x_{1}\right|+\left|x_{2}\right|+\cdots+\left|x_{d}\right|\right)^{p}=\sum_{|\alpha| \leq p} \frac{p!}{(p-|\alpha|)!\alpha!}\left|x^{\alpha}\right|
$$

Hence for $|\beta| \leq p$, one finds

$$
\begin{aligned}
(1+|x|)^{p}\left|\partial^{\beta} K\right| & \leq \frac{1}{(2 \pi)^{d}} \sum_{|\alpha| \leq p} \frac{p!}{(p-|\alpha|)!\alpha!}\left|\partial^{\alpha} \phi_{\beta}\right|_{L^{1}\left(\mathbb{R}^{d}\right)} e^{|\operatorname{Im} x|} \\
& \leq \frac{1}{(2 \pi)^{d}}\left(\sum_{|\alpha| \leq p} \frac{p!}{(p-|\alpha|)!\alpha!} \sup _{|\beta| \leq p}\left\{\left|\partial^{\alpha} \phi_{\beta}\right|_{L^{1}\left(\mathbb{R}^{d}\right)}\right\}\right) e^{|\operatorname{Im} x|}
\end{aligned}
$$

Thus

$$
\left|\partial^{\beta} K(x)\right| \leq c_{p} \frac{e^{|\operatorname{Im} x|}}{(1+|x|)^{p}}
$$

proving (115). Moreover, if $x \in \mathbb{R}^{d}$ and $\beta \in \mathbb{N}^{d}$, we have:

$$
\begin{aligned}
\left|\partial^{\beta} K(x)\right| & \leq \frac{1}{(2 \pi)^{d}} \int_{\mathbb{R}^{d}}\left|u^{\beta} \phi(u)\right| d u \\
& \leq \frac{1}{(2 \pi)^{d}}\|\phi\|_{L^{1}}
\end{aligned}
$$

which is (116).
The relation (115) shows that the integral $I_{\alpha, \beta}$ is well defined. Moreover, notice that if $\alpha_{i}>\beta_{i}$ for some $i, \alpha_{i}$ integrations by part show that $I_{\alpha, \beta}=0$ since $\partial^{\alpha_{i}} u^{\beta}=0$. So we can assume that $\alpha_{i} \leq \beta_{i}$ for all $i$ (i.e. $\alpha \leq \beta$ ). By integration by part, we see that

$$
I_{\alpha, \beta}=(-1)^{|\alpha|} \frac{\beta!}{(\beta-\alpha)!} \int_{\mathbb{R}^{d}} u^{\beta-\alpha} K(u+i v) d u
$$

Notice also that

$$
\begin{aligned}
K(u+i v) & =\frac{1}{(2 \pi)^{d}} \int_{\mathbb{R}^{d}} \phi(\xi) e^{i \xi \cdot(u+i v)} d \xi \\
& =\frac{1}{(2 \pi)^{d}} \int_{\mathbb{R}^{d}} \phi(\xi) e^{-v \cdot \xi} e^{i \xi \cdot u} d \xi=: \tilde{K}_{v}(u)
\end{aligned}
$$

that we can think to be the anti-Fourier transform of

$$
\tilde{\phi}_{v}(\xi)=\phi(\xi) e^{-v \cdot \xi} \in C_{0}^{\infty}\left(\mathbb{R}^{d}\right)
$$

for each fixed $v$.
Hence:

$$
\begin{aligned}
\int_{\mathbb{R}^{d}} u^{\beta-\alpha} K(u+i v) d u & =\left.\frac{1}{(-i)^{|\beta-\alpha|}} \frac{\partial^{|\beta-\alpha|}}{\partial \xi^{\beta-\alpha}}\right|_{\xi=0} \int_{\mathbb{R}^{d}} K(u+i v) e^{-i u \xi} d u \\
& =\left.\frac{1}{(-i)^{|\beta-\alpha|}} \frac{\partial^{|\beta-\alpha|}}{\partial \xi^{\beta-\alpha}}\right|_{\xi=0} \int_{\mathbb{R}^{d}} \tilde{K}_{v}(u) e^{-i u \xi} d u \\
& =\left.\frac{1}{(-i)^{|\beta-\alpha|}} \frac{\partial^{|\beta-\alpha|}}{\partial \xi^{\beta-\alpha}}\right|_{\xi=0}\left(\phi(\xi) e^{-v \xi}\right)
\end{aligned}
$$

where, in the last equalities, we used the fact that the right hand-side integral is the Fourier transform of $\tilde{K}_{v}$ (i.e. $\tilde{\phi}_{v}$ ). Now, using (8), we obtain:

$$
\begin{aligned}
\int_{\mathbb{R}^{d}} u^{\beta-\alpha} K(u+i v) d u & =\frac{1}{(-i)^{|\beta-\alpha|}} \phi(0)(-v)^{\beta-\alpha} \\
& =(-i v)^{\beta-\alpha}
\end{aligned}
$$

From the properties above, we can conclude:

$$
I_{\alpha, \beta}:=\int_{\mathbb{R}^{d}} u^{\beta} \partial^{\alpha} K(u+i v) d u=\left\{\begin{array}{lr}
(-1)^{|\beta|} \frac{\beta!}{(\beta-\alpha)!}(i v)^{\beta-\alpha}, & \text { if } \alpha \leq \beta \\
0, & \text { otherwise }
\end{array}\right.
$$

We proceed with the proof of Proposition 8. Recall the definition of $f_{r}=\mathcal{S}_{r}$ given in (114) and notice that (117) implies $\mathcal{S}_{r} P=P$ for any polynomial $P$ :

$$
\begin{aligned}
\left(\mathcal{S}_{r} P\right)(x) & =\int_{\mathbb{R}^{d}} K(\eta) P(x-r \eta) d \eta \\
& :=\int_{\mathbb{R}^{d}} K(\eta)\left(\sum_{|k| \leq \operatorname{deg} P} a_{k}(x, r) \eta^{k}\right) d \eta \\
& =\sum_{|k| \leq \operatorname{deg} P} a_{k}(x, r) \int_{\mathbb{R}^{d}} K(\eta) \eta^{k} d \eta \\
& =a_{0}(x, r)=P(x)
\end{aligned}
$$

We claim also that for every $l \in \mathbb{R}_{+}$and $f \in C^{l}$, then there exists a constant $c_{0}(l, d)$ such that

$$
\begin{equation*}
\left|f(x+y)-P_{k}(x, y)\right| \leq c_{0}(l, d)|f|_{C^{l}}|y|^{l} \tag{118}
\end{equation*}
$$

where

$$
P_{k}(x, y)=\sum_{|\alpha| \leq k} \frac{\partial^{\alpha} f(x)}{\alpha!} y^{\alpha} \quad \text { and } \quad k=[l] .
$$

In fact, if $l$ is an integer (118) comes immediately from Taylor's formula (and, actually, one could get $c_{0}(l, d)$ as small as one wants); if $l=k+\mu$ with $\mu \in(0,1)$, we get

$$
\begin{aligned}
\left|f(x+y)-P_{k}(x, y)\right| & \leq\left|\int_{0}^{1} \frac{(1-t)^{k-1}}{(k-1)!}\left(\partial^{k} f(x+t y) y^{(k)}-\partial^{k} f(x) y^{(k)}\right) d t\right| \\
& \leq \int_{0}^{1} \frac{(1-t)^{k-1}}{(k-1)!} \frac{\left|\partial^{k} f(x+t y)-\partial^{k} f(x)\right|}{|t y|^{\mu}}|y|^{k+\mu} d t \\
& \leq \frac{|f|_{C^{l}}}{k!}|y|^{l}
\end{aligned}
$$

Let $x=u+i v, \eta=\frac{u-y}{r}$ where $u, v, y$ belong to $\mathbb{R}^{d}$. Then

$$
\begin{aligned}
\partial^{\alpha} f_{r}(x) & =\partial^{\alpha}\left(\frac{1}{r^{d}} \int_{\mathbb{R}^{d}} K\left(\frac{x-y}{r}\right) f(y) d y\right) \\
& =\frac{1}{r^{d}} \int_{\mathbb{R}^{d}} \frac{1}{r^{|\alpha|}} \partial^{\alpha} K\left(\frac{u-y}{r}+i \frac{v}{r}\right) f(y) d y \\
& =\frac{1}{r^{|\alpha|}} \int_{\mathbb{R}^{d}} \frac{1}{r^{d}} \partial^{\alpha} K\left(\eta+i \frac{v}{r}\right) f(u-r \eta) d(u-r \eta) \\
& =\frac{1}{r^{|\alpha|}} \int_{\mathbb{R}^{d}} \partial^{\alpha} K\left(\eta+i \frac{v}{r}\right) f(u-r \eta) d \eta .
\end{aligned}
$$

Let us consider now $|\beta| \leq l-|\alpha|$; in view of (117):

$$
\begin{aligned}
\frac{\partial^{\alpha+\beta} f(u)}{\beta!}(i v)^{\beta} & =\frac{\partial^{\alpha+\beta} f(u)}{\beta!} \frac{(i v)^{\beta}}{(-1)^{|\alpha|} \frac{(\beta+\alpha)!}{\beta!}\left(-i \frac{v}{r}\right)^{\beta}} \int_{\mathbb{R}^{d}} \partial^{\alpha} K\left(\eta+i \frac{v}{r}\right) \eta^{\alpha+\beta} d \eta \\
& =\int_{\mathbb{R}^{d}} \partial^{\alpha} K\left(\eta+i \frac{v}{r}\right) \frac{\partial^{\alpha+\beta} f(u)}{(\beta+\alpha)!}(-1)^{|\alpha+\beta|} r^{|\beta|} \eta^{\alpha+\beta} d \eta \\
& =\frac{1}{r^{|\alpha|}} \int_{\mathbb{R}^{d}} \partial^{\alpha} K\left(\eta+i \frac{v}{r}\right) \frac{\partial^{\alpha+\beta} f(u)}{(\beta+\alpha)!}(-1)^{|\alpha+\beta|}(r \eta)^{\alpha+\beta} d \eta
\end{aligned}
$$

Hence, if we denote $k=[l]$ and apply again (117), we obtain:

$$
\begin{aligned}
& \sum_{|\beta| \leq l-|\alpha|} \frac{\partial^{\alpha+\beta} f(u)}{\beta!}(i v)^{\beta} \\
& \quad=\frac{1}{r^{|\alpha|}} \int_{\mathbb{R}^{d}} \partial^{\alpha} K\left(\eta+i \frac{v}{r}\right) \sum_{|\beta| \leq l-|\alpha|}\left(\frac{\partial^{\alpha+\beta} f(u)}{(\beta+\alpha)!}(-1)^{|\alpha+\beta|}(r \eta)^{\alpha+\beta}\right) d \eta \\
& \quad=\frac{1}{r^{|\alpha|}} \int_{\mathbb{R}^{d}} \partial^{\alpha} K\left(\eta+i \frac{v}{r}\right) P_{k}(u-r \eta) d \eta
\end{aligned}
$$

In view of the above calculations and of the inequality (118), we have:

$$
\begin{aligned}
\mid \partial^{\alpha} f_{r}(x) & \left.-\sum_{|\beta| \leq l-|\alpha|} \partial^{\alpha+\beta} f(u) \frac{(i v)^{\beta}}{\beta!} \right\rvert\, \\
& \leq \frac{1}{r^{|\alpha|}} \int_{\mathbb{R}^{d}}\left|\partial^{\alpha} K\left(\eta+i \frac{v}{r}\right)\right|\left|f(u-r \eta)-P_{k}(u-r \eta)\right| d \eta
\end{aligned}
$$

$$
\begin{aligned}
& \leq c_{0}|f|_{C^{l}} r^{l-|\alpha|} \int_{\mathbb{R}^{d}}\left|\partial^{\alpha} K\left(\eta+i \frac{v}{r}\right)\right||\eta|^{l} d \eta \\
& \leq c_{0}|f|_{C^{l}} r^{l-|\alpha|} \int_{\mathbb{R}^{d}}\left|\partial^{\alpha} K\left(\eta+i \frac{v}{r}\right)\right|(1+|\eta|)^{l} d \eta .
\end{aligned}
$$

Applying (117) for $p>l$ big enough (i.e., $p>l+d$ ), we get

$$
\begin{aligned}
\left|\partial^{\alpha} f_{r}(x)-\sum_{|\beta| \leq l-|\alpha|} \partial^{\alpha+\beta} f(u) \frac{(i v)^{\beta}}{\beta!}\right| & \leq c_{1}|f|_{C^{l}} r^{l-|\alpha|} \int_{R^{d}}(1+|\eta|)^{l-p} d \eta \\
& \leq c(l, d)|f|_{C^{l}} r^{l-|\alpha|}
\end{aligned}
$$

which completes the proof.

Remark 15 (i) In particular, (113) (with $x \in \mathbb{R}^{d}$ ) implies, for $s \leq l(s$ integer),

$$
\begin{equation*}
\left|f_{r}-f\right|_{C^{s}} \leq c^{\prime}|f|_{C^{l}} r^{l-s}, \quad(s \leq l) \tag{119}
\end{equation*}
$$

for a suitable $c^{\prime}=c^{\prime}(l, d)$.
(ii) Also, (113) with $l=0$, yields, for every $f \in C^{0}$ and any $r>0$,

$$
\begin{equation*}
\sup _{\Delta_{r}^{d}}\left|f_{r}\right| \leq(c+1)|f|_{C^{0}} \tag{120}
\end{equation*}
$$

(iii) If $x$ is real, the definition of $f_{r}$ implies immediately that, if $f \in C^{l}$ and $|\alpha| \leq l$,

$$
\begin{equation*}
\partial^{\alpha} f_{r}(x)=\left(\partial^{\alpha} f\right)_{r}(x) . \tag{121}
\end{equation*}
$$

By analyticity, such relation is seen to hold for any $x \in \mathbb{C}^{d}$.
(iv) Using the observation in (iii) and Cauchy estimates, one can give the following bound on the derivatives of $f_{r}$ with $f \in C^{l}$ :

$$
\begin{equation*}
\left|f_{r}\right|_{C^{s}} \leq c^{\prime \prime}|f|_{C^{l}} r^{l-s}, \quad(s \geq l \text { integers }) \tag{122}
\end{equation*}
$$

$c^{\prime \prime}=c^{\prime \prime}(l, d)$ being a suitable positive constant.
Proof If $s=l$, (122) comes from (119). Let $s>l$. For any multi-index $\alpha$ such that $|\alpha|=s$ we can find $\beta$ and $\alpha_{0}$ such that $\alpha=\beta+\alpha_{0}$ with $\left|\alpha_{0}\right|=l$ and $|\beta|=s-l$. Then,
by (121), Cauchy estimates and (120), denoting by " const " (possibly different) constants depending on $l$ and $d$, we find

$$
\begin{aligned}
\left|\partial^{\alpha} f_{r}\right|_{C^{0}} & =\left|\partial^{\beta}\left(\partial^{\alpha_{0}} f\right)_{r}\right|_{C^{0}} \leq \text { const } \sup _{\Delta_{r}^{d}}\left|\left(\partial^{\alpha_{0}} f\right)_{r}\right| \frac{1}{r^{s-l}} \\
& \leq \text { const }\left|\partial^{\alpha_{0}} f\right|_{C^{0}} \frac{1}{r^{s-l}} \leq \text { const }|f|_{C^{l}} r^{l-s}
\end{aligned}
$$

(v) (Convexity estimates) Let $l>0$, let $f \in C^{l}$ and let $k, m$ be integers such that $0 \leq k \leq m \leq l$. Then, there exist a constant $\hat{c}=\hat{c}(l, d)>0$ such that

$$
\begin{equation*}
|f|_{C^{m}} \leq \hat{c}|f|_{C^{k}}^{\frac{l-m}{l-k}}|f|_{C^{l}}^{\frac{m-k}{l-k}} \tag{123}
\end{equation*}
$$

Proof Define $r:=\left(\frac{|f|_{C_{k}}}{|f|_{C^{l}} l}\right)^{\frac{1}{l-k}}$. Then $r \leq 1$ and by (119) and (122), we get
$|f|_{C^{m}} \leq\left|f_{r}-f\right|_{C^{m}}+\left|f_{r}\right|_{C^{m}} \leq \mathrm{const}\left(|f|_{C^{l}} r^{l-m}+|f|_{C^{k}} r^{k-m}\right)=2$ const $|f|_{C^{k}}^{\frac{l-m}{T-k}}|f|_{C^{l}}^{\frac{m-k}{l-k}}$.
(vi) If $f$ is periodic, $f_{r}$ is obviously periodic. Moreover, if $f$ belongs to $C^{0}\left(\mathbb{T}^{d}\right)$, then

$$
\begin{aligned}
\left(\mathcal{S}_{r} f\right)_{n} & =f_{\mathbb{T}^{d}} \mathcal{S}_{r} f(x) e^{-i x \cdot n} d x \\
& =f_{\mathbb{T}^{d}}\left(\int_{\mathbb{R}^{d}} K(\eta) f(x-r \eta) d \eta\right) e^{-i x \cdot n} d x \\
& =\int_{\mathbb{R}^{d}} K(\eta) e^{-i r n \cdot \eta}\left(f_{\mathbb{T}^{d}} f(x-r \eta) e^{-i(x-r \eta) \cdot n} d x\right) d \eta \\
& =f_{n} \phi(r n)
\end{aligned}
$$

showing that $f_{r}$ is a trigonometric polynomial.

Proposition 10 (Bernstein, Moser) Let $l>0$ and $d \in \mathbb{Z}_{+}$. Let $f_{0}=0$ and for each $j$ in $\mathbb{Z}_{+}$, let $f_{j}$ be a real analytic function on $\Delta_{r_{j}}^{d} \subset \mathbb{C}^{d}$ where $r_{j}=r_{0} / 2^{j}$ for some $0<r_{0} \leq 1$. Assume that

$$
\begin{equation*}
\left|f_{j}-f_{j-1}\right|_{r_{j}} \leq A r_{j}^{l} \tag{124}
\end{equation*}
$$

for every $j \geq 1$ and some constant $A$.
Then, $f_{j}$ tends to $f$ uniformly on $\mathbb{R}^{d}$ and $f \in C^{s}\left(\mathbb{R}^{d}\right)$ for every non integer
$s \leq l$. Furthermore, there exists a constant $C=C(l, d)$ such that:

$$
\begin{equation*}
|f|_{C^{s}\left(\mathbb{R}^{d}\right)} \leq \frac{C A}{\mu(1-\mu)} r_{0}^{l-s} \tag{125}
\end{equation*}
$$

where $\mu=s-[s]$. Finally, if the $f_{i}$ 's are periodic in each variable $x_{j}$ then so is $f$.

Remark 16 (i) If $f_{0} \neq 0$ and (124) holds for all $j \geq 1$, we can apply the proposition to $\tilde{f}_{j}:=f_{j}-f_{0}$ getting that $\tilde{f}_{j}$ tends uniformly to $\tilde{f} \in C^{s}$ so that $f_{j}$ tends uniformly to $f:=f_{0}+\tilde{f} \in C^{s}$. Moreover,

$$
\begin{equation*}
\left|f-f_{0}\right|_{C^{s}\left(\mathbb{R}^{d}\right)} \leq \frac{C A}{\mu(1-\mu)} r_{0}^{l-s} . \tag{126}
\end{equation*}
$$

(ii) It is enough to prove Proposition 10 in the particular case where $l \in(0,1)$ and $s=l=\mu$ as we proceed to check.

Proof of point (ii) of Remark 16 Let us consider the three following claims:
(a) "Proposition 10 holds true for $0<s<l=1$ ".
(b) "Proposition 10 holds true for $0<s<l<1$ ".
(c) "Proposition 10 holds true for $0<s=l<1$ ".

We will show that

$$
(c) \Rightarrow(b) \Rightarrow(a)
$$

and finally that (a) implies the general case.
$(b) \Rightarrow(a)$ : To prove (a), we assume (124) with $l=1$ and fix $0<s<1$. Then, for every $s<l<1$,

$$
\left|f_{j}-f_{j-1}\right|_{r_{j}} \leq A r_{j} \leq A r_{j}^{l}
$$

which shows that $f_{j}$ satisfies the hypothesis (124) of (b). Applying (b), we get

$$
|f|_{C^{s}} \leq \frac{C A}{s(1-s)} r_{0}^{l-s}
$$

and taking the infimum in the above expression over $l<1$, we get

$$
|f|_{C^{s}} \leq \frac{C A}{s(1-s)} r_{0}^{1-s}
$$

$(c) \Rightarrow(b)$ : We have

$$
\left|f_{j}-f_{j-1}\right|_{r_{j}} \leq A r_{j}^{l}=A r_{j}^{l-s} r_{j}^{s} \leq A r_{0}^{l-s} r_{j}^{s}
$$

which shows that $\tilde{f}_{j}=f_{j} / r_{0}^{l-s}$ satisfies the hypothesis of $(c)$. Then, by (c), the uniform limit $\tilde{f}$ of $\tilde{f}_{j}$ belongs to $C^{s}$ and

$$
|\tilde{f}|_{C^{s}} \leq \frac{C A}{s(1-s)}
$$

which is equivalent to

$$
|f|_{C^{s}} \leq \frac{C A}{s(1-s)} r_{0}^{l-s}
$$

which proves (b).
Now, let us show that the claim (c) implies the general case. We prove by induction on $k$ ( $k \geq 1$ ) that
$\left(\mathcal{P}_{k}\right) \quad$ "Proposition 10 holds true for $0<l \leq k$ ".
First of all, notice that $\left(\mathcal{P}_{1}\right)$ holds true since $(c)$ implies $(a)$ and $(b)$. Let $s$ be a non integer such that $0<s \leq l \leq k+1$. We can assume that $k<l \leq k+1$ (if not, then $0<l \leq k$ and we can apply the inductive hypothesis). By assumption, we have

$$
\left|f_{j}-f_{j-1}\right|_{r_{j}} \leq A r_{j}^{l} .
$$

Using Cauchy estimates (Lemma 5) we have for every $\alpha \in \mathbb{N}^{d}$ such that $|\alpha|=1$,

$$
\left|\partial^{\alpha} f_{j}-\partial^{\alpha} f_{j-1}\right| \frac{r_{j}}{2} \leq\left|f_{j}-f_{j-1}\right|_{r_{j}}\left(\frac{r_{j}}{2}\right)^{-1} \leq 2 A r_{j}^{l-1}
$$

Then by $\left(\mathcal{P}_{k}\right), f_{j}$ converges uniformly to $f \in C^{s}$ for any $s \leq l-1$ and

$$
|f|_{C^{s}} \leq \frac{C(l-1) A}{\mu(1-\mu)} r_{0}^{l-1-s} .
$$

Proof of Proposition 10 In view of point (ii) of Remark 16, we may suppose, without loss of generality, that

$$
\begin{equation*}
0<l=s=\mu<1 \tag{127}
\end{equation*}
$$

In this case, we have to prove that $f_{j}$ converges uniformly on $\mathbb{R}^{d}$ to $f$ and that

$$
|f|_{C^{\mu}}=|f|_{C^{0}}+\sup _{0<|x-y| \leq 1} \frac{|f(x)-f(y)|}{|x-y|^{\mu}} \leq \frac{C A}{\mu(1-\mu)}
$$

Set $g_{j}=f_{j}-f_{j-1}$. First of all, let us prove that $f_{j}$ converges uniformly on $\mathbb{R}^{d}$. For any $1 \leq n \leq N$, one has (recall that $r_{j}=r_{0} / 2^{j}$ )

$$
\begin{align*}
\left|\sum_{j=n}^{N} g_{j}\right|_{C^{0}} & \leq \sum_{j=n}^{N}\left|g_{j}\right|_{C^{0}}=\sum_{j=n}^{N}\left|f_{j}-f_{j-1}\right|_{C^{0}}  \tag{128}\\
& =\sum_{j=n}^{\infty} A\left(\frac{r_{0}}{2^{j}}\right)^{\mu}=\frac{A r_{0}^{\mu}}{2^{\mu n}} \frac{1}{1-2^{-\mu}}
\end{align*}
$$

which converges to zero as $n$ goes to $+\infty$; thus $f_{j}$ converges to $f=\sum_{j=1}^{\infty} g_{j}$ uniformly on $\mathbb{R}^{d}$ and (setting $n=1$ and $N=+\infty$ in the above estimates) we have

$$
|f|_{C^{0}} \leq A r_{0}^{\mu} \frac{2^{-\mu}}{1-2^{-\mu}}
$$

Since $1-2^{-\mu} \geq \mu / 2$ for each $\mu \in[0,1]$, we get:

$$
\begin{equation*}
|f|_{C^{0}} \leq \frac{2 A r_{0}^{\mu}}{\mu}<\frac{2 A}{\mu(1-\mu)} r_{0}^{\mu} \leq \frac{2 A}{\mu(1-\mu)} \tag{129}
\end{equation*}
$$

In order to estimate the second part of $|f|_{C^{\mu}}$ we need to distinguish two cases according to whether $r_{0}<|x-y| \leq 1$ or $|x-y| \leq r_{0}$.

First case: $r_{0}<|x-y| \leq 1$. Then using the second inequality in (129), we get

$$
\begin{equation*}
|f(x)-f(y)| \leq 2|f|_{C^{0}} \leq \frac{4 A}{\mu(1-\mu)} r_{0}^{\mu} \leq \frac{4 A}{\mu(1-\mu)}|x-y|^{\mu} \tag{130}
\end{equation*}
$$

Second case: $0<|x-y| \leq r_{0}$. Then there exists a $N$ in $\mathbb{N}$ such that:

$$
\begin{equation*}
\frac{r_{0}}{2^{N+1}} \leq|x-y| \leq \frac{r_{0}}{2^{N}} \tag{131}
\end{equation*}
$$

The second inequality in (129) is equivalent to

$$
\begin{equation*}
\left(2^{N} r_{0}^{-1}\right)^{1-\mu} \leq|x-y|^{\mu-1} \tag{132}
\end{equation*}
$$

Now,

$$
\begin{aligned}
|f(x)-f(y)| & \leq \sum_{j=1}^{\infty}\left|g_{j}(x)-g_{j}(y)\right| \\
& =\sum_{j=1}^{N}\left|g_{j}(x)-g_{j}(y)\right|+\sum_{j=N+1}^{\infty}\left|g_{j}(x)-g_{j}(y)\right|
\end{aligned}
$$

Let us estimate separately the two sums. Using Cauchy estimates, we have (recall that, by hypothesis, $\left|g_{j}\right|_{r_{j}} \leq A r_{j}^{\mu}$ ):

$$
\left|\partial_{x} g_{j}\right| \frac{r_{j}}{2} \leq\left|g_{j}\right|_{r_{j}}\left(\frac{r_{j}}{2}\right)^{-1} \leq\left|g_{j}\right|_{C^{0}}\left(\frac{r_{j}}{2}\right)^{-1} \leq 2 A r_{j}^{\mu-1}
$$

Hence,

$$
\begin{aligned}
\sum_{j=1}^{N}\left|g_{j}(x)-g_{j}(y)\right| & \leq 2 A|x-y| \sum_{j=1}^{N}\left(\frac{2^{j}}{r_{0}}\right)^{1-\mu} \\
& =2 A|x-y|\left(\left(r_{0}^{-1}\right)^{1-\mu}\right) \frac{2^{(N+1)(1-\mu)}-1}{2^{1-\mu}-1}
\end{aligned}
$$

Since $2^{t}-1 \geq t / 2$ for any $t \geq 0$ (and since $2^{1-\mu} \leq 2$ ), by (132) we get

$$
\begin{align*}
\sum_{j=1}^{N}\left|g_{j}(x)-g_{j}(y)\right| & \leq 2 A|x-y|\left(\left(r_{0}^{-1}\right)^{1-\mu}\right)\left(4 \frac{2^{(N+1)(1-\mu)}}{1-\mu}\right) \\
& \leq 16 A \frac{|x-y|^{\mu}}{1-\mu} \tag{133}
\end{align*}
$$

Next (using again $1-2^{-\mu} \geq \mu / 2$ for $\mu \in[0,1]$ and (129))

$$
\begin{align*}
\sum_{j=N+1}^{\infty}\left|g_{j}(x)-g_{j}(y)\right| & \leq 2 \sum_{j=N+1}^{\infty}\left|g_{j}\right|_{C^{0}} \leq 2 A \sum_{j=N+1}^{\infty}\left(\frac{r_{0}}{2^{j}}\right)^{\mu} \\
& \leq 2 A\left(\frac{r_{0}}{2^{N+1}}\right)^{\mu} \frac{1}{1-2^{-\mu}} \leq 4 A \frac{|x-y|^{\mu}}{\mu} \tag{134}
\end{align*}
$$

Putting (133) and (134) together, we get:

$$
\begin{align*}
|f(x)-f(y)| & \leq \frac{16 A}{1-\mu}|x-y|^{\mu}+\frac{4 A}{\mu}|x-y|^{\mu} \\
& \leq \frac{16 A}{\mu(1-\mu)}|x-y|^{\mu} \tag{135}
\end{align*}
$$

Thus, by (129) and (135), we get

$$
|f|_{C^{\mu}} \leq \frac{C A}{\mu(1-\mu)}
$$

with $C=18$.

### 2.2 A KAM theorem in $C^{k}$ category

In this section we extend KAM theory to the finitely differentiable case.
For simplicity, we shall discuss only the nearly integrable case: in particular we prove the following generalization of Corollary 1.

Theorem 11 Let $\omega \in \mathbb{R}^{d}$ be $(\gamma, \tau)$-diophantine, let $l>l_{0}:=4 \tau+3$, let $V \in C^{l}\left(\mathbb{T}^{d}\right)$ and let $M>0$ be such that $|V|_{C^{l}} \leq M$. There exists a constant $\kappa=\kappa(l, d, \tau, \gamma, M)>1$ such that if

$$
\begin{equation*}
\kappa\left(\left|V_{x}\right|_{C^{0}}\right)^{\frac{l-l_{0}}{l-1}} \leq 1 \tag{136}
\end{equation*}
$$

then there exists a function $u: \mathbb{T}^{d} \rightarrow \mathbb{R}^{d}$, which belongs to $C^{s}$ for all $s \leq l-l_{0}$ not integer, satisfying

$$
\begin{equation*}
D^{2} u+V_{x}(\theta+u)=0 \tag{137}
\end{equation*}
$$

and

$$
\begin{equation*}
|u|_{C^{s}} \leq \frac{\kappa}{\mu(1-\mu)}\left(\left|V_{x}\right|_{C^{0}}\right)^{\frac{l-l_{0}-s}{l-1}}, \quad \mu:=s-[s] \tag{138}
\end{equation*}
$$

If $s<2$, from the proof given below it follows easily (as relation (137) suggests) that the double directional derivative $D^{2} u$ exists and is a $C^{s}\left(\mathbb{T}^{d}\right)$ function (exercise).

Proof Let

$$
\begin{equation*}
\varepsilon:=\left(\left|V_{x}\right|_{C^{0}}\right)^{\frac{1}{l-1}}, \quad \xi_{j}:=\frac{\varepsilon}{2^{j}}, \quad \hat{\xi}_{j}:=\frac{\xi_{j+1}}{2}=\frac{\xi_{j}}{4} . \tag{139}
\end{equation*}
$$

Notice that (136) implies that $\varepsilon<1$.
By Proposition 8, the real-analytic functions $V_{j}:=\mathcal{S}_{\xi_{j}} V \in \mathcal{R}_{\xi_{j}}$ satisfy

$$
\begin{equation*}
\left|\partial^{\alpha} V_{j}(x)-\sum_{|\beta| \leq l-|\alpha|} \frac{\partial^{\beta+\alpha} V(\operatorname{Re} x)}{\beta!}(i \operatorname{Im} x)^{\beta}\right| \leq c|V|_{C^{l}} \xi_{j}^{l-|\alpha|}, \tag{140}
\end{equation*}
$$

for every $x \in \Delta_{\xi_{j}}^{d}$ and $|\alpha| \leq l$. Denote by $\mathcal{E}_{j}$ the differential operator

$$
\mathcal{E}_{j}: v \rightarrow \mathcal{E}_{j}(v):=D^{2} v+\partial_{x} V_{j}(\theta+v) .
$$

The strategy is to construct a sequence of real-analytic functions $u_{j} \in \mathcal{R}_{\hat{\xi}_{j}}$, satisfying $\mathcal{E}_{j}\left(u_{j}\right)=0$ and to obtain, by Proposition 10 , the solution $u$ as uniform limit of the $u_{j}^{\prime} s$.

For the purpose of this proof, we denote by "const " (possibly different) constants depending on $l, d$ and $\tau$ and by $\kappa_{i}$ suitable constants depending on $l, d, \tau, \gamma, M$. The constant $\kappa$ in (136) is assumed to be such that

$$
\begin{equation*}
\kappa \geq \kappa_{i}, \quad \forall i \tag{141}
\end{equation*}
$$

As a preliminary remark, we observe that, for any $|\alpha| \leq 3$ and for any $j \geq 0$,

$$
\begin{equation*}
\sup _{\Delta_{\xi_{j}}^{d}}\left|\partial^{\alpha} V_{j}\right| \leq \text { const } M \tag{142}
\end{equation*}
$$

as it follows from (140) and the fact that $\varepsilon<1$ :

$$
\begin{aligned}
\left|\partial^{\alpha} V_{j}(x)\right| \leq & \left|\partial^{\alpha} V_{j}(x)-\sum_{|\beta| \leq l-3} \frac{\partial^{\beta+\alpha} V(\operatorname{Re} x)}{\beta!}(i \operatorname{Im} x)^{\beta}\right| \\
& +\left|\sum_{|\beta| \leq l-3} \frac{\partial^{\beta+\alpha} V(\operatorname{Re} x)}{\beta!}(i \operatorname{Im} x)^{\beta}\right| \\
\leq & \text { const }\left(M \xi_{j}^{l-3}+|V|_{C^{3}}+M \xi_{j}\right) \\
& \leq \text { const } M .
\end{aligned}
$$

We proceed in three steps: construction of $u_{0}$; inductive construction of $u_{j}$ $(j \geq 1)$; construction of $u$ as $\lim u_{j}$.

Step 1: construction of $u_{0}$. We want to apply the KAM Theorem 6 with $v \equiv 0, \xi_{*}=\xi_{0}, \xi=\xi_{0} / 2:=\varepsilon / 2, \bar{\xi}=\hat{\xi}_{0}=\varepsilon / 4$. We start by estimating $\mathcal{E}_{0}(v)=\mathcal{E}_{0}(0)$. Let $\theta \in \Delta_{\xi}^{d}$, i.e., $\left|\operatorname{Im} \theta_{k}\right| \leq \varepsilon / 2$. Then

$$
\begin{aligned}
\left|\mathcal{E}_{0}(0)(\theta)\right| & :=\left|\partial_{x} V_{0}(\theta)\right| \\
\leq & \left|\partial_{x} V_{0}(\theta)-\sum_{|\beta| \leq l-1} \frac{\partial^{\beta} \partial_{x} V(\operatorname{Re} \theta)}{\beta!}(i \operatorname{Im} \theta)^{\beta}\right| \\
& +\left|\sum_{|\beta| \leq l-1} \frac{\partial^{\beta} \partial_{x} V(\operatorname{Re} \theta)}{\beta!}(i \operatorname{Im} \theta)^{\beta}\right|
\end{aligned}
$$

$$
\left.\begin{array}{l}
\stackrel{(*)}{\leq} \\
\stackrel{(* *)}{(*)} \\
\stackrel{c o n s t}{\leq}\left(|V|_{C^{l}} \varepsilon^{l-1}+\sum_{j \leq l-1}\left|V_{x}\right|_{C^{j}} \varepsilon^{j}\right) \\
\stackrel{(\dagger)}{=} \quad \operatorname{const}\left(|V|_{C^{l}} \varepsilon^{l-1}+\sum_{j \leq l-1}\left|V_{x^{\prime}}\right|_{C^{0}}^{\frac{l-1-j}{l-1}}|V|_{C^{l}}^{\frac{j}{l-1}} \varepsilon^{j}\right) \\
\leq  \tag{143}\\
= \\
= \\
=
\end{array} \kappa_{1} \varepsilon^{l-1}, \quad \text { const }\left.\max \{1, M\}\right|_{C^{l}}+\sum_{j \leq l-1}|V|_{C^{l}}^{\frac{j}{l-1}}\right) \varepsilon^{l-1} .
$$

where: $(*)$ is implied by (140); $(* *)$ is the convexity estimate (123) with $k=0, m=j$ and $l$ replaced by $l-1 ;(\dagger)$ is the definition of $\varepsilon$. Thus, recalling the notations in Theorem 6 , and observing that $\alpha \leq \max \left\{1\right.$, const $\left.M / \gamma^{2}\right\}$ by (142), that $\lambda=\eta=1$ and recalling the definition of $\hat{\xi}_{j}$, we see that (75) is implied, in our case, by

$$
\begin{equation*}
\kappa_{2} \varepsilon^{l-l_{0}} \leq 1 \tag{144}
\end{equation*}
$$

for a suitable $\kappa_{2}>1$. Such condition, in view of (141) and of the definition of $\varepsilon$, is implied by (136). Therefore, by Theorem 6 , there exists a function $u_{0} \in \mathcal{R}_{\hat{\xi}_{0}}$ such that

$$
\mathcal{E}_{0}\left(u_{0}\right)=0
$$

and such that

$$
\begin{equation*}
\left\|u_{0}\right\|_{\hat{\xi}_{0}},\left\|\partial_{\theta} u_{0}\right\|_{\hat{\xi}_{0}} \leq \kappa_{3} \varepsilon^{l-l_{0}} \leq 1 \tag{145}
\end{equation*}
$$

where $\kappa_{3}=K \kappa_{2}, K$ being the constant in (77); the second inequality holds because of (141) and (136). The first step is completed.

Step 2: construction of $\left\{u_{j}\right\}$. We proceed inductively constructing $u_{j+1}$, for $j \geq 0$, via Theorem 6 by taking $v=u_{j}$ as approximate solution. We also take $\xi_{*}=\xi_{j+1}, \xi=\hat{\xi}_{j}$ and $\bar{\xi}=\hat{\xi}_{j+1}$. The parameter $\alpha$, in view of (142), is uniformly bounded by $\max \left\{1\right.$, const $\left.M / \gamma^{2}\right\}$.

We, now, assume that, for $0 \leq k \leq j$, there exist functions $u_{k} \in \mathcal{R}_{\hat{\xi}_{k}}$ such that

$$
\mathcal{E}_{k}\left(u_{k}\right)=0
$$

and such that ${ }^{22}$

$$
\begin{equation*}
\left\|u_{k}-u_{k-1}\right\|_{\hat{\xi}_{k}},\left\|\partial_{\theta}\left(u_{k}-u_{k-1}\right)\right\|_{\hat{\xi}_{k}} \leq \kappa_{4}\left(\frac{\varepsilon}{2^{k}}\right)^{l-l_{0}}, \quad(1 \leq k \leq j) \tag{146}
\end{equation*}
$$

[^15]for a suitable $\kappa_{4} \geq \kappa_{3}$ specified below; finally we assume that $0 \leq k \leq j$ :
\[

$$
\begin{equation*}
\left\|\partial_{\theta} u_{k}\right\|_{\hat{\xi}_{k}} \leq 1, \quad\left\|\left(I+\partial_{\theta} u_{k}\right)^{-1}\right\|_{\hat{\xi}_{k}} \leq 2 . \tag{147}
\end{equation*}
$$

\]

Notice that $\left\|\partial_{\theta} u_{k}\right\|_{\hat{\xi}_{k}} \leq 1$ implies $^{23}\left\|I+\partial_{\theta} u_{k}\right\|_{\hat{\xi}_{k}} \leq 2$ so that, if (147) holds, then in Theorem 6 one can take $\lambda=\eta=2$. The inductive assumption (147) and the definitions in (139) imply that (74) is satisfied: in fact, if $\theta \in \Delta_{\hat{\xi}_{j}}^{d}$, then

$$
\begin{aligned}
\left|\operatorname{Im} u_{j}(\theta)\right| & =\left|\operatorname{Im}\left(u_{j}(\theta)-u_{j}(\operatorname{Re} \theta)\right)\right| \\
& \leq\left|u_{j}(\theta)-u_{j}(\operatorname{Re} \theta)\right| \\
& \leq\left\|\partial_{\theta} u_{j}\right\|_{\hat{\xi}_{j}} \hat{\xi}_{j} \\
& \leq \hat{\xi}_{j}:=\xi_{j+1}-\hat{\xi}_{j}
\end{aligned}
$$

We need, now, to estimate $\mathcal{E}_{j+1}\left(u_{j}\right)$. Since, by the inductive assumption, $\mathcal{E}_{j}\left(u_{j}\right)=0$, we find, for $\theta \in \Delta_{\hat{\xi}_{j}}^{d}$ and because of (140),

$$
\begin{align*}
&\left|\mathcal{E}_{j+1}\left(u_{j}\right)(\theta)\right|:=\left|D^{2} u_{j}(\theta)+\partial_{x} V_{j+1}\left(\theta+u_{j}\right)\right| \\
&=\left|\partial_{x} V_{j+1}\left(\theta+u_{j}\right)-\partial_{x} V_{j}\left(\theta+u_{j}\right)\right| \\
& \quad \leq\left|\partial_{x} V_{j+1}\left(\theta+u_{j}\right)-\sum_{|\beta| \leq l-1} \frac{\partial_{x}^{\beta} \partial_{x} V\left(\operatorname{Re}\left(\theta+u_{j}\right)\right)}{\beta!}\left(i \operatorname{Im}\left(\theta+u_{j}\right)\right)^{\beta}\right| \\
&+\left|\sum_{|\beta| \leq l-1} \frac{\partial_{x}^{\beta} \partial_{x} V\left(\operatorname{Re}\left(\theta+u_{j}\right)\right)}{\beta!}\left(i \operatorname{Im}\left(\theta+u_{j}\right)\right)^{\beta}-\partial_{x} V_{j}\left(\theta+u_{j}\right)\right| \\
& \leq \quad \text { const }|V|_{C^{l}}\left(\xi_{j+1}^{l-1}+\xi_{j}^{l-1}\right) \mid \\
& \leq \text { const } M \xi_{j}^{l-1} . \tag{148}
\end{align*}
$$

Thus, (75) becomes, in the present case,

$$
\begin{equation*}
\kappa_{5}\left(\frac{\varepsilon}{2^{j+1}}\right)^{l-l_{0}} \leq 1 \tag{149}
\end{equation*}
$$

for a suitable $\kappa_{5}>1$. We now define ${ }^{24} \kappa_{4}$ as

$$
\begin{equation*}
\kappa_{4}:=\max \left\{\kappa_{3}, K \kappa_{5}\right\} \tag{150}
\end{equation*}
$$

[^16]Notice that condition (149) is again implied by (136). Thus, by Theorem 6, there exists a function $u_{j+1} \in \mathcal{R}_{\hat{\xi}_{j+1}}$ such that

$$
\mathcal{E}_{j+1}\left(u_{j+1}\right)=0
$$

and such that

$$
\begin{equation*}
\left\|u_{j+1}-u_{j}\right\|_{\hat{\xi}_{j+1}},\left\|\partial_{\theta}\left(u_{j+1}-u_{j}\right)\right\|_{\hat{\xi}_{j+1}} \leq \kappa_{4}\left(\frac{\varepsilon}{2^{j+1}}\right)^{l-l_{0}} \tag{151}
\end{equation*}
$$

which is exactly (146) with $k=j+1$. The bounds (151) together with the condition (136) easily implies that the inductive assumptions (147) are satisfied also for ${ }^{25} k=j+1$, allowing to iterate the inductive procedure indefinitely. The second step is completed.

Step 3: construction of $u$. At this point we can apply Proposition 10 (see also Remark 16) with: $l$ replaced by $l-l_{0} ; f_{j}=u_{j}-u_{0} ; r_{j}=\hat{\xi}_{j}:=\varepsilon / 2^{j+2}$ (so that $r_{0}=\varepsilon / 4$ ); $A=\kappa_{4} 4^{l-2}$ (compare (124), (151) and the choice of $r_{j}$ ). The thesis of the theorem now follows at once from Proposition 10.

Exercise Discuss the $C^{\infty}$ case.
Exercise* Extend Theorem 6 to the differentiable case.

[^17]
## 3 Appendix A

Lemma 12 If $T=T(\theta)$ is a strictly positive and symmetric real matrix for each $\theta \in \mathbb{T}^{d}$, then

$$
\|\langle T\rangle\|^{-1} \leq \sup _{\theta \in \mathbb{T}^{d}}\left\|T^{-1}\right\|
$$

Proof By hypotheses there exists an orthogonal matrix $P$ such that $P^{T} T P$ is diagonal. Let $\left\{\lambda_{i}: i=1, \ldots, d\right\}$ be the spectrum of $T$ and $y$ a vector whose coordinates are $y_{i}$ for $i=1, \ldots, d$ in the basis where $T$ is diagonal. Then, we have:

$$
\begin{aligned}
(T y) \cdot y & =\sum_{i=1}^{d} \lambda_{i} y_{i}^{2} \\
& \geq \min _{i \in\{1, \ldots, d\}}\left\{\lambda_{i}\right\}\|y\|^{2} .
\end{aligned}
$$

But, $\min _{i \in\{1, \ldots, d\}}\left\{\lambda_{i}\right\}=\left\|T^{-1}\right\|^{-1}$, thus,

$$
\begin{aligned}
(T y) \cdot y & \geq \frac{\|y\|^{2}}{\left\|T^{-1}\right\|} \\
& \geq \frac{\|y\|^{2}}{\sup _{\theta \in \mathbb{T}^{d}}\left\|T^{-1}\right\|}
\end{aligned}
$$

Set $y=\langle T\rangle^{-1} x$. Taking the average of the last expression, we get

$$
\left\langle\left(T\langle T\rangle^{-1} x\right) \cdot y\right\rangle \geq\left\langle\frac{\|y\|^{2}}{\sup _{\theta \in \mathbb{T}^{d}}\left\|T^{-1}\right\|}\right\rangle
$$

i.e.,

$$
\left(\langle T\rangle\langle T\rangle^{-1} x\right) \cdot y \geq \frac{\|y\|^{2}}{\sup _{\theta \in \mathbb{T}^{d}}\left\|T^{-1}\right\|}
$$

Finally, using Schwarz's inequality in the left-hand side, and dividing by $\|y\|$, we get

$$
\|x\| \geq \frac{\|y\|}{\sup _{\theta \in \mathbb{T}^{d}}\left\|T^{-1}\right\|}
$$

i.e.,

$$
\sup _{\theta \in \mathbb{T}^{d}}\left\|T^{-1}\right\| \geq \frac{\|y\|}{\|x\|}=\frac{\left\|\langle T\rangle^{-1} x\right\|}{\|x\|}
$$

Since this is true for all $x$, the result follows.

Proof of Lemma 4 To prove (66), we observe that by Parseval identity, if $v \in \mathbb{C}^{d}$ is such that $|\operatorname{Im} v|<\xi$, then

$$
\begin{equation*}
\sum_{n \in \mathbb{Z}^{d}}\left|f_{n}\right|^{2} e^{-2(n \cdot v)}=\int_{\mathbb{T}^{d}}|f(u+i v)|^{2} d u \leq\|f\|_{L^{2}, \Delta_{\xi}^{d}}^{2} \tag{152}
\end{equation*}
$$

Thus, for any $n \in \mathbb{Z}^{d}$ and any $v$ as above,

$$
e^{-n \cdot v}\left|f_{n}\right| \leq\|f\|_{L^{2}, \Delta_{\xi}^{d}}
$$

and choosing ${ }^{26} v=-\left(\operatorname{sign} n_{1}, \ldots, \operatorname{sign} n_{d}\right)(\xi-\varepsilon)$ we get

$$
e^{|n| \xi}\left|f_{n}\right| \leq\|f\|_{L^{2}, \Delta_{\xi}^{d}} e^{\varepsilon|n|}
$$

letting $\varepsilon \rightarrow 0$, (66) follows.
Let us turn to (67). For the purpose of the following argument we let $|n|$ denote the Euclidean norm also for integer vectors.
The first inequality in (67) is obvious; in order to establish the second one for a fixed vector $x \in \Delta_{\xi-\delta}^{d}$ we define the set

$$
I_{0}=\left\{n \in \mathbb{Z}^{d}: \quad(n \cdot \operatorname{Im} x) \leq-\frac{|n|(\xi-\delta)}{2}\right\}
$$

and let

$$
\mu:=\frac{\xi}{\xi-\delta}
$$

Observe that $(n \cdot \operatorname{Im} x) \leq-\frac{|n|(\xi-\delta)}{2}$ is equivalent to

$$
\begin{equation*}
-n \cdot \operatorname{Im} x \leq-n \cdot \mu \operatorname{Im} x-\frac{|n| \delta}{2} \tag{153}
\end{equation*}
$$

[^18]Thus, by (153), Schwarz inequality and (152) (since $\mu x \in \Delta_{\xi}^{d}$ ), we get

$$
\begin{aligned}
\sum_{n \in I_{0}}\left|f_{n}\right| e^{-(n \cdot \operatorname{Im} x)} & \leq \sum_{n \in I_{0}}\left|f_{n}\right| e^{-(n \cdot \mu \operatorname{Im} x)} e^{-|n| \delta} \\
& \leq\left(\sum_{n \in \mathbb{Z}^{d}}\left|f_{n}\right|^{2} e^{-2(n \cdot \mu \operatorname{Im} x)}\right)^{\frac{1}{2}}\left(\sum_{n \in \mathbb{Z}^{d}} e^{-|n| \delta}\right)^{\frac{1}{2}} \\
& \leq \frac{c_{1}^{\frac{1}{2}}}{\delta^{\frac{d}{2}}}\|f\|_{L^{2}, \Delta_{\xi}^{d}}
\end{aligned}
$$

where ${ }^{27}$

$$
c_{1}=c_{1}(d)=\sup _{0<\lambda \leq 1} \sum_{n \in \mathbb{Z}^{d}} \lambda^{d} e^{-|n| \lambda}<\infty .
$$

It is easy to see that there exist an integer $s=s(d)$ and a collection of $s$ unit vectors $e_{1}, \ldots, e_{s}$ in $\mathbb{R}^{d}$, such that for ever $y \in \mathbb{R}^{d}$ there exists $\sigma \in\{1, \ldots, s\}$ with ${ }^{28}$

$$
\left(y \cdot e_{\sigma}\right)>\frac{|y|}{2} .
$$

Now, every nonzero integer vector outside $I_{0}$, lies in one the sets $I_{\sigma}$ defined as

$$
I_{\sigma}:=\left\{n \in \mathbb{Z}^{d}: \quad(n \cdot \operatorname{Im} x)>-\frac{|n|(\xi-\delta)}{2}, \quad\left(n \cdot e_{\sigma}\right)>\frac{|n|}{2}\right\} .
$$

But, (using again Schwarz inequality and (152)),

$$
\begin{aligned}
\sum_{n \in I_{\sigma}}\left|f_{n}\right| e^{-(n \cdot \operatorname{Im} x)} & \leq \sum_{n \in I_{\sigma}}\left|f_{n}\right| e^{|n| \frac{(\xi-\delta)}{2}} \\
& \leq\left(\sum_{n \in I_{\sigma}}\left|f_{n}\right|^{2} e^{|n| \xi}\right)^{\frac{1}{2}}\left(\sum_{n \in I_{\sigma}} e^{-|n| \delta}\right)^{\frac{1}{2}} \\
\sum e^{-|n| \delta}=\frac{1}{\delta^{d}}\left(\delta^{d}\right. & \left.\sum e^{-|n| \delta}\right) \leq \frac{1}{\delta^{d}} \sup _{0<\lambda<1}\left(\lambda^{d} \sum e^{-|n| \lambda}\right) .
\end{aligned}
$$

${ }^{28}$ Given $v \in S^{d-1}:=\left\{y \in \mathbb{R}^{d}:|y|=1\right\}$, let $C_{v}:=\left\{w \in S^{d-1}: w \cdot v>\frac{1}{2}\right\}$. Then $C_{v}$ is an open (in the relative topology) neighborhood of $v$ and $\left\{C_{v}: v \in S^{d-1}\right\}$ is an open cover of the compact set $S^{d-1}$. Thus there exist unit vectors $v_{1}=: e_{1}, \ldots, v_{s}=: e_{s}$ such that $S^{d-1} \subset \cup_{j=1}^{s} C_{e_{j}}$ : this is equivalent to the claim.

$$
\begin{aligned}
& \leq\left(\sum_{n \in \mathbb{Z}^{d}}\left|f_{n}\right|^{2} e^{-2\left(n \cdot \xi e_{\sigma}\right)}\right)^{\frac{1}{2}}\left(\sum_{n \in \mathbb{Z}^{d}} e^{-|n| \delta}\right)^{\frac{1}{2}} \\
& \leq \frac{c_{1}^{\frac{1}{2}}}{\delta^{\frac{d}{2}}}\|f\|_{L^{2}, \Delta_{\xi}^{d}}
\end{aligned}
$$

Inequality (67) now follows and one can take $c_{0}=(s+1) \sqrt{c_{1}}$.

## 4 Appendix B (Fourier Norms)

For $\xi \geq 0$ let us define the space

$$
\hat{\mathcal{R}}_{\xi}\left(\mathbb{T}^{d}, \mathbb{R}^{N}\right):=\left\{f \in C\left(\mathbb{T}^{d}, \mathbb{R}^{N}\right) \text { s. t. }\|f\|_{\xi}^{\hat{}}:=\sum_{n \in \mathbb{Z}^{d}}\left|f_{n}\right| e^{|n| \xi}<\infty\right\}
$$

The space $\hat{\mathcal{R}}_{\xi}\left(\mathbb{T}^{d}, \mathbb{R}^{N}\right)$ is a Banach space with respect to the norm $\|f\|_{\hat{\xi}}$. Moreover, since

$$
\left|f_{n}\right| \leq\|f\|_{\xi} e^{-|n| \xi}
$$

(when $\xi>0$ ) the function $f$ has a holomorphic extension to the complex strip $\Delta_{\xi}^{d}$. Notice that

$$
\begin{equation*}
\sup _{\Delta_{\xi}^{d}}|f|=\sup _{\Delta_{\xi}^{d}}\left|\sum_{n \in \mathbb{Z}^{d}} f_{n} e^{i n x}\right| \leq \sum_{n \in \mathbb{Z}^{d}}\left|f_{n}\right| e^{|n| \xi}=\|f\|_{\xi}^{\hat{}} \tag{154}
\end{equation*}
$$

and if $0 \leq \xi^{\prime}<\xi$, then $\hat{\mathcal{R}}_{\xi} \varsubsetneqq \hat{\mathcal{R}}_{\xi^{\prime}}$ (exercise). In particular, (154) shows that $f \in \hat{\mathcal{R}}_{\xi}$ admits a holomorphic and bounded extension to $\Delta_{\xi}^{d}$.

Lemma 13 Let $f \in \hat{\mathcal{R}}_{\xi}\left(\mathbb{T}^{d}, X\right), p \in \mathbb{Z}, \alpha \in \mathbb{N}^{d}$ be such that $|p|+|\alpha|>0$. If $p>0$, assume either $|\alpha|>0$ or $\langle f\rangle=0$. Let $0<\delta \leq \xi$. Then

$$
\left\|D^{-p} \partial^{\alpha} f\right\|_{\xi-\delta} \leq C_{p, \alpha}(\omega)\|f\|_{\xi}
$$

where

$$
C_{p, \alpha}(\omega):=\sup _{n \neq 0} \frac{\left|n^{\alpha}\right| e^{-\delta|n|}}{|\omega \cdot n|^{p}} .
$$

If $\omega$ is $(\gamma, \tau)$-diophantine, then

$$
C_{p, \alpha}(\omega) \leq \begin{cases}\frac{(p \tau+|\alpha|)!}{\gamma^{p} \delta^{p \tau+|\alpha|}}, & \text { for } p \geq 0 \\ \left(\sup \left|\omega_{i}\right|\right)^{|p|} \frac{(|p|+|\alpha|)!}{\delta^{|p|+|\alpha|}}, & \text { for } p<0\end{cases}
$$

Proof

$$
\left\|D^{-p} \partial^{\alpha} f\right\|_{\xi-\delta}=\left\|\sum_{n \neq 0} \frac{n^{\alpha} f_{n}}{(\omega \cdot n)^{p}} e^{i n x}\right\|_{\xi-\delta}^{\hat{n}}=\sum_{n \neq 0} \frac{\left|f_{n}\right|\left|n^{\alpha}\right|}{|\omega \cdot n|^{p}} e^{|n|(\xi-\delta)} \leq C_{p, \alpha}(\omega)\left\|f_{n}\right\|_{\xi}
$$

If $p \geq 0$ and $\omega$ is $(\gamma, \tau)$-diophantine, then

$$
C_{p, \alpha}(\omega) \leq \sup _{n \neq 0} \frac{|n|^{|\alpha|+\tau p}}{\gamma^{p}} e^{-|n| \delta}
$$

The function on the right is of the form $g(t)=t^{a} e^{-\delta|t|}, t>0, a \geq 0$, and has a maximum at the point $t_{m}=\frac{a}{\delta}$ such that $g\left(t_{m}\right)=a^{a}(e \delta)^{-a} \leq a!\delta^{-a}$, where if $a$ is not an integer, we define ${ }^{29} a!=([a]+1)$ !
If now $p<0$, one can repeat the previous arguments using the fact that in this case $|\omega \cdot n|^{|p|} \leq\left(\sup \left|\omega_{i}\right|\right)^{|p|}|n|^{|p|}$.

In particular, if $\langle f\rangle=0$, then

$$
\begin{equation*}
\left\|D^{-1} f\right\|_{\xi-\delta} \leq\|f\|_{\xi} \hat{\frac{\tau!}{\gamma \delta^{\tau}}} \tag{155}
\end{equation*}
$$

Lemma 14 Let $f \in \hat{\mathcal{R}}_{\xi}\left(\mathbb{T}^{d}, X\right), g \in \hat{\mathcal{R}}_{\xi}\left(\mathbb{T}^{d}, Y\right)$, with $X$ and $Y$ tensor spaces $\left(\mathbb{R}^{N}\right.$, matrices or higher dimensional tensors) and assume that the product fg is well defined. Then

$$
\|f g\|_{\hat{\xi}} \leq\|f\|_{\hat{\xi}}\|g\|_{\xi}^{\hat{k}}
$$

${ }^{29}$ If $a \in \mathbb{N}$, then $\left(\frac{a}{e}\right)^{a}=\frac{a^{a}}{1+a+\frac{a^{2}}{2}+\ldots+\frac{a^{a}}{a!}+\ldots} \leq \frac{a^{a}}{\frac{a^{a}}{a!}}=a!$. If $a$ is non-integer, then we can repeat the same argument eliminating all terms in the Taylor expansion for the exponent except $1+a^{[a]+1} /([a]+1)!$.

Proof Indeed,

$$
\begin{aligned}
\|f g\|_{\xi} & =\sum_{n}\left|(f g)_{n}\right| e^{|n| \xi}=\sum_{n}\left|\sum_{m} f_{n-m} g_{m}\right| e^{|n| \xi} \\
& \leq \sum_{n, m}\left|f_{n-m}\left\|g_{m} \mid e^{|n-m| \xi} e^{|m| \xi}=\right\| f\left\|_{\hat{\xi}}\right\| g \|_{\xi}^{\hat{\xi}}\right.
\end{aligned}
$$

Lemma 15 Let $V \in \hat{\mathcal{R}}_{\xi_{*}}\left(\mathbb{T}^{d}, \mathbb{R}\right), g \in \hat{\mathcal{R}}_{\xi}\left(\mathbb{T}^{d}, \mathbb{R}^{d}\right)$, with $\xi_{*}>\xi$. If $\left\|g_{i}\right\|_{\xi} \leq$ $\xi_{*}-\xi$ for all $i=1, \ldots, d$, then $\|V(x+g(x))\|_{\dot{\xi}} \leq\|V\|_{\xi_{*}}$.

Proof Using Lemma 14 , the fact that $\|\cdot\|_{\xi}$ is a norm, one finds

$$
\begin{aligned}
\|V(x+g(x))\|_{\xi} & =\sum_{n} \mid\left(V(x+g(x))_{n} \mid e^{|n| \xi}\right. \\
& =\sum_{n}\left|\left(\sum_{m} V_{m} e^{i m \cdot(x+g(x))}\right)_{n}\right| e^{|n| \xi} \\
& =\sum_{n}\left|\sum_{m} V_{m}\left(e^{i m \cdot g(x)}\right)_{n-m}\right| e^{|n| \xi} \\
& \leq \sum_{n, m}\left|V_{m}\right|\left|\left(e^{i m \cdot g(x)}\right)_{n-m}\right| e^{|n| \xi} \\
& \left.\leq \sum_{n, m}\left|V_{m}\right| \sum_{j \geq 0}\left[\frac{(i m \cdot g(x))^{j}}{j!}\right]_{n-m} \right\rvert\, e^{|n-m| \xi} e^{|m| \xi} \\
& \leq \sum_{m, j} \frac{\left|V_{m}\right|}{j!} e^{|m| \xi} \sum_{n}\left|\left[(i m \cdot g(x))^{j}\right]_{n-m}\right| e^{|n-m| \xi} \\
& \leq \sum_{m, j} \frac{\left|V_{m}\right|}{j!} e^{|m| \xi}\left(\|m \cdot g\|_{\xi}\right)^{j} \\
& \leq \sum_{m, j} \frac{\left|V_{m}\right|}{j!} e^{|m| \xi}\left(\sum_{i=1}^{d} m_{i}\left\|g_{i} \mid\right\|_{\xi}\right)^{j} \\
& \leq \sum_{m, j} \frac{\left|V_{m}\right|}{j!} e^{|m| \xi}\left(\sup _{i}\left\|g_{i}\right\|_{\xi}^{\wedge}\right)|m|^{j} \\
& \leq \sum_{m, j} \frac{\left|V_{m}\right|}{j!} e^{|m| \xi}\left(\xi_{*}-\xi\right)^{j}|m|^{j} \\
& =\sum_{m}\left|V_{m}\right| e^{|m| \xi} e^{\left(\xi_{*}-\xi\right)|m|}=\|V\|_{\xi_{*}} .
\end{aligned}
$$

The Fourier norm $\|\cdot\|_{\xi}$ and the sup-norm $\|\cdot\|_{\xi}$ are not equivalent (exercise); however they are strictly related. We have already seen (compare (154)) that

$$
\begin{equation*}
\|f\|_{\xi} \leq\|f\|_{\xi}, \tag{156}
\end{equation*}
$$

which implies immediately ${ }^{30}$

$$
\begin{equation*}
\hat{\mathcal{R}}_{\xi}\left(\mathbb{T}^{d}\right) \subset \mathcal{R}_{\xi}\left(\mathbb{T}^{d}\right) \tag{157}
\end{equation*}
$$

We now prove that a weaker version of the converse of (156) is true. Let $\xi^{\prime}>\xi>0$ and assume that $f \in \mathcal{R}_{\xi^{\prime}}$. Since, for every $n \in \mathbb{Z}^{d}$,

$$
\left|f_{n}\right| \leq\|f\|_{\xi^{\prime}} e^{-|n| \xi^{\prime}}
$$

we have

$$
\|f\|_{\xi}=\sum_{n}\left|f_{n}\right| e^{|n| \xi} \leq\|f\|_{\xi^{\prime}} \sum_{n} e^{-|n|\left(\xi^{\prime}-\xi\right)} .
$$

But (for suitable positive constants $c(d), C(d)$ )

$$
\sum_{n \in \mathbb{Z}^{d}} e^{-|n|\left(\xi^{\prime}-\xi\right)} \leq c(d) \int_{\mathbb{R}^{d}} e^{-|x|\left(\xi^{\prime}-\xi\right)} d x=\frac{c(d)}{\left(\xi^{\prime}-\xi\right)^{d}} \int_{\mathbb{R}^{d}} e^{-|y|} d y=\frac{C(d)}{\left(\xi^{\prime}-\xi\right)^{d}}
$$

so that

$$
\begin{equation*}
\|f\|_{\xi} \leq C(d) \frac{\|f\|_{\xi^{\prime}}}{\left(\xi^{\prime}-\xi\right)^{d}} \tag{158}
\end{equation*}
$$

This relation shows, in particular, that

$$
\begin{equation*}
\mathcal{R}_{\xi^{\prime}}\left(\mathbb{T}^{d}\right) \subset \hat{\mathcal{R}}_{\xi}\left(\mathbb{T}^{d}\right), \quad \forall \xi^{\prime}>\xi \tag{159}
\end{equation*}
$$

Exercise Give explicit upper bounds on $c(d)$ and $C(d)$.

[^19]
[^0]:    *Published in "Dynamical Systems. Part I: Hamiltonian Systems and Celestial Mechanics", Pubblicazioni della Classe di Scienze, Scuola Normale Superiore, Pisa. 1-56. Centro di Ricerca Matematica "Ennio De Giorgi" : Proceedings. (2003)

[^1]:    ${ }^{1}$ I.e., quasi periodic-solutions with $d$ independent frequencies; for the definition of quasiperiodic solutions, see below.
    ${ }^{2}$ Such technique gives precise hypotheses in order to approximate $C^{l}$ functions with realanalytic ones and, viceversa, to get $C^{l}$ functions out of limits of real-analytic sequences.

[^2]:    3 "." denotes, here, the standard inner product.

[^3]:    ${ }^{4}$ I.e., the equation $F_{y}\left(u_{0}, x\right) w_{0}+\varepsilon_{0}=0$.

[^4]:    ${ }^{5}$ For vectors $a, b \in \mathbb{R}^{d}$ we denote $a \cdot b:=\sum_{i=1}^{d} a_{i} b_{i}$
    ${ }^{6}$ In fact, if $\omega T=2 \pi n$ with $T>0$ and $n \in \mathbb{Z}^{d} \backslash\{0\}$ then there are exactly $(d-1)$ independent vectors $n_{j}$ s.t. $\omega \cdot n_{j}=0$.

[^5]:    ${ }^{7}$ As well known, $\theta \in \mathbb{T}^{d} \rightarrow \omega t \in \mathbb{T}^{d}$ is dense if and only if $\omega$ is rationally independent; see, e.g., [V.I. Arnold, Mathematical methods of classical mechanics, Springer-Verlag, 1989].

[^6]:    ${ }^{9}$ At an intuitive level, one should think to substitute the error function $\varepsilon$ with $\mu \varepsilon$ thinking $\mu$ as a small real parameter: the terms appearing with a $\mu$ in front will be thought of as "small" terms and of the same "order" of the error function, terms with a $\mu^{2}$ in front will be thought of as "quadratically smaller terms", etc.
    ${ }^{10}$ Here and in what follows the symbol $Q_{j}$ 's stand for terms "quadratic in $\varepsilon$ ".

[^7]:    ${ }^{11} \mathcal{R}_{0}$ denotes simply $C\left(\mathbb{T}^{d}, \mathbb{R}^{N}\right)$ endowed with the sup-norm.
    ${ }^{12}$ For example, if $X=\operatorname{Mat}(n \times n)$, then $f_{n} \in X$ and, in the definition of the norm $\|f\|_{\xi}$, the expression $\left|f_{n}\right|$ denotes the standard "operator norm" $\sup _{|c|=1}\left|f_{n} c\right|$.

[^8]:    ${ }^{13}$ Exercise.

[^9]:    ${ }^{14}$ Notice that $f$ is bounded on $\Delta_{\xi}^{d}$ so $f$ cannot have singularities on the boundary of the strip.

[^10]:    ${ }^{15}$ Recall Definition 2.
    ${ }^{16}$ Recall that $\mathcal{E}$ is the differential operator defined as $\mathcal{E}(v):=D^{2} v+V_{x}(\theta+v)$ where $D:=\sum_{i=1}^{d} \omega_{i} \partial_{\theta_{i}}$; notice, also, that, since $\left\|\operatorname{Im} v_{j}\right\|_{\xi} \leq \xi_{*}-\xi$, then $\theta+v \in \Delta_{\xi_{*}}^{d}$ whenever $\theta \in \Delta_{\xi}^{d}$.

[^11]:    ${ }^{17}$ Later we shall make a specific (somewhat arbitrary) choice.
    ${ }^{18}$ For the proof see Lemma 12 in Appendix A.

[^12]:    ${ }^{19}$ This follows from the following inequality: $(1-x)^{-1} \leq 1+2 x$ valid for $0 \leq x \leq 1 / 2$.

[^13]:    ${ }^{20}$ We will use the following notations:

    $$
    \operatorname{Re} x:=\left(\operatorname{Re} x_{1}, \ldots, \operatorname{Re} x_{d}\right) \quad \text { and } \quad \operatorname{Im} x:=\left(\operatorname{Im} x_{1}, \ldots, \operatorname{Im} x_{d}\right)
    $$

[^14]:    ${ }^{21}$ For vectors $\alpha, \beta \in \mathbb{N}^{d}$, we denote

    $$
    \alpha \leq \beta \quad \Longleftrightarrow \quad \alpha_{i} \leq \beta_{i} \quad \forall i=1, \ldots, d
    $$

[^15]:    ${ }^{22}$ For $j=0$ (146) is obviously replaced by the already proven (145).

[^16]:    ${ }^{23}$ We are choosing norms for which $\|I\|=1$.
    ${ }^{24}$ This is well defined since in the computations leading to the definition of $\kappa_{5}$ the inductive hypotheses (146) have not been used.

[^17]:    ${ }^{25}$ Exercise Fill in the details.

[^18]:    ${ }^{26}$ Here, we let signa be 1 if $a \geq 0$ and $(-1)$ otherwise.

[^19]:    ${ }^{30}$ Actually $\hat{\mathcal{R}}_{\xi}\left(\mathbb{T}^{d}\right) \varsubsetneqq \mathcal{R}_{\xi}\left(\mathbb{T}^{d}\right)$ (exercise).

