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# Reducibility of the Quantum Harmonic Oscillator in $d$-dimensions with Polynomial Time Dependent Perturbation 

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#### Abstract

We prove a reducibility result for a quantum harmonic oscillator in arbitrary dimension with arbitrary frequencies perturbed by a linear operator which is a polynomial of degree two in $x_{j},-\mathrm{i} \partial_{j}$ with coefficients which depend quasiperiodically on time.


## 1 Introduction and statement

The aim of this paper is to present a reducibility result for the time dependent Schrödinger equation

$$
\begin{array}{r}
\mathrm{i} \dot{\psi}=H_{\epsilon}(\omega t) \psi, x \in \mathbb{R}^{d} \\
H_{\epsilon}(\omega t):=H_{0}+\epsilon W(\omega t, x,-\mathrm{i} \nabla) \tag{1.2}
\end{array}
$$

where

$$
\begin{equation*}
H_{0}:=-\Delta+V(x), \quad V(x):=\sum_{j=1}^{d} \nu_{j}^{2} x_{j}^{2}, \quad \nu_{j}>0 \tag{1.3}
\end{equation*}
$$

and $W(\theta, x, \xi)$ is a real polynomial in $(x, \xi)$ of degree at most two, with coefficients being real analytic functions of $\theta \in \mathbb{T}^{n}$. Here $\omega$ are parameters which are assumed to belong to the set $\mathcal{D}=(0,2 \pi)^{n}$.

For $\epsilon=0$ the spectrum of 1.2 is given by

$$
\begin{equation*}
\sigma\left(H_{0}\right)=\left\{\lambda_{k}\right\}_{k \in \mathbb{N}^{d}}, \quad \lambda_{k} \equiv \lambda_{\left(k_{1}, \ldots, k_{d}\right)}:=\sum_{j=1}^{d}\left(2 k_{j}+1\right) \nu_{j} \tag{1.4}
\end{equation*}
$$

[^0]with $k_{j} \geq 0$ integers. In particular if the frequencies $\nu_{j}$ are nonresonant, then the differences between couples of eigenvalues are dense on the real axis. As a consequence, in the case $\epsilon=0$ most of the solutions of 1.1 are almost periodic with an infinite number of rationally independent frequencies.

Here we will prove that for any choice of the mechanical frequencies $\nu_{j}$ and for $\omega$ belonging to a set of large measure in $\mathcal{D}$ the system 1.1 is reducible: precisely there exists a time quasiperiodic unitary transformation of $L^{2}\left(\mathbb{R}^{d}\right)$ which conjugates 1.2 ) to a time independent operator; we also deduce boundedness of the Sobolev norms of the solution.

The proof exploits the fact that for polynomial Hamiltonians of degree at most 2 the correspondence between classical and quantum mechanics is exact (i.e. without error term), so that the result can be proven by exact quantization of the classical KAM theory which ensures reducibility of the classical Hamiltonian system

$$
\begin{equation*}
h_{\epsilon}:=h_{0}+\epsilon W(\omega t, x, \xi), \quad h_{0}:=\sum_{j=1}^{d} \xi_{j}^{2}+\nu_{j}^{2} x_{j}^{2} . \tag{1.5}
\end{equation*}
$$

We will use (in the appendix) the exact correspondence between classical and quantum dynamics of quadratic Hamiltonians also to prove a complementary result. Precisely we will present a class of examples (following GY00]) in which one generically has growth of Sobolev norms. This happens when the frequencies $\omega$ of the external forcing are resonant with some of the $\nu_{j}$ 's.

We recall that the exact correspondence between classical and quantum dynamics of quadratic Hamiltonians was already exploited in the paper HLS86 to prove stability/instability results for one degree of freedom time dependent quadratic Hamiltonians.

Notwithstanding the simplicity of the proof, we think that the present result could have some interest, since this is the first example of a reducibility result for a system in which the gaps of the unperturbed spectrum are dense in $\mathbb{R}$. Furthermore it is one of the few cases in which reducibility is obtained for systems in more than one space dimension.

Indeed, most of the results on the reducibility problem for 1.1) have been obtained in the one dimensional case, and also the results in higher dimensions obtained up to now deal only with cases in which the spectrum of the unperturbed system has gaps whose size is bounded from below, like in the Harmonic oscillator (or in the Schrödinger equation on $\mathbb{T}^{d}$ ). On the other hand we restrict here to perturbations, which although unbounded, must belong to the very special class of polynomials in $x_{j}$ and $-\mathrm{i} \partial_{j}$. The reason is that for operators in this class, the commutator is the operator whose symbol is the Poisson bracket of the corresponding symbols, without any error term (see Remark 2.2 and Remark 2.4. In order to deal with more general perturbations one needs further ideas and techniques.

Before closing this introduction we recall some previous works on the reducibility problem for (1.1) and more generally for perturbations of the Schrödinger equation with a potential $V(x)$. As we already anticipated, most of the works deal with the one dimensional case. The first one is Com87 in which pure point nature of the Floquet operator is obtained in case of a smoothing perturbation of the Harmonic oscillator in dimension 1 (see also Kuk93). The techniques of this paper were extended in DS96, DLSV02], in order to deal with potentials growing superquadratically (still in dimension 1 ) but with perturbations which were only required to be bounded.

A slightly different approach originates from the so called KAM theory for PDEs Kuk87, Way90. In particular the methods developed in that context in order to deal with unbounded perturbations (see Kuk97, Kuk98]) where exploited in BG01 in order to deal with the reducibility problem of 1.1) with superquadratic potential in dimension 1 (see [LY10] for a further
improvement). The case of bounded perturbations of the Harmonic oscillator in dimension 1 was treated in Wan08, GT11.

An extension of KAM theory to NLS on $\mathbb{T}^{d}$ has been obtained in EK10 and its methods have been adapted to deal with the reducibility problem of quasiperiodically forced linear Schrödinger equation in EK09. A further reducibility result for equations in more then one space dimensions is GP16, in which bounded perturbations of the completely resonant Harmonic Oscillator in $\mathbb{R}^{d}$ were studied. As far as we know, these are the only higher dimensional linear systems for which reducibility is known ${ }^{11}$.

We remark that all these papers deal with cases where the spectrum of the unperturbed operator is formed by well separated eigenvalues. In the higher dimensional cases they are allowed to have high multiplicity localized in clusters. But then the perturbation must have special properties ensuring that the clusters are essentially not destroyed under the KAM iteration.

Finally we recall the works Bam17a, Bam17b in which pseudodifferential calculus was used together with KAM theory in order to prove reducibility results for (in dimension 1) with unbounded perturbations. The ideas of the present paper are a direct development of the ideas of Bam17a, Bam17b. We also recall that the idea of using pseudodifferential calculus together with KAM theory in order to deal with problems involving unbounded perturbations originates from the work [PT01, IPT05] and has been developed in order to give a quite general theory in [BBM14, BM16, Mon14] (see also [FP15]).

In order to state our main result, we need some preparations. It is well known that the equation (1.1) is well posed (see for example MR17) in the scale $\mathcal{H}^{s}, s \in \mathbb{R}$ of the weighted Sobolev spaces defined as follows. For $s \geq 0$ let

$$
\mathcal{H}^{s}:=\left\{\psi \in L^{2}\left(\mathbb{R}^{d}\right): \quad H_{0}^{s / 2} \psi \in L^{2}\left(\mathbb{R}^{d}\right)\right\}
$$

equipped with the natural Hilbert space norm $\|\psi\|_{s}:=\left\|H_{0}^{s / 2} \psi\right\|_{L^{2}\left(\mathbb{R}^{d}\right)}$. For $s<0, \mathcal{H}^{s}$ is defined by duality. Such spaces are not dependent on $\nu$ for $\nu_{j}>0,1 \leq j \leq d$. We also have $\mathcal{H}^{s} \equiv \operatorname{Dom}\left(-\Delta+|x|^{2}\right)^{s / 2}$.

We will prove the following reducibility theorem:
Theorem 1.1. Let $\psi$ be a solution of (1.1). There exist $\epsilon_{*}>0, C>0$ and $\forall|\epsilon|<\epsilon_{*}$ a closed set $\mathcal{E}_{\epsilon} \subset(0,2 \pi)^{n}$ with $\operatorname{meas}\left((0,2 \pi)^{n} \backslash \mathcal{E}_{\epsilon}\right) \leq C \epsilon^{\frac{1}{9}}$ and, $\forall \omega \in \mathcal{E}_{\epsilon}$ there exists a unitary (in $L^{2}$ ) time quasiperiodic map $U_{\omega}(\omega t)$ s.t. defining $\varphi$ by $U_{\omega}(\omega t) \varphi=\psi$, it satisfies the equation

$$
\begin{equation*}
\mathrm{i} \dot{\varphi}=H_{\infty} \varphi, \tag{1.6}
\end{equation*}
$$

with $H_{\infty}$ a positive definite time independent operator which is unitary equivalent to a diagonal operator

$$
\sum_{j=1}^{d} \nu_{j}^{\infty}\left(x_{j}^{2}-\partial_{x_{j}}^{2}\right)
$$

where $\nu_{j}^{\infty}=\nu_{j}^{\infty}(\omega)$ are defined for $\omega \in \mathcal{E}_{\epsilon}$ and fulfill the estimates

$$
\left|\nu_{j}-\nu_{j}^{\infty}\right| \leq C \epsilon, \quad j=1, \ldots, d
$$

Finally the following properties hold

[^1](i) $\forall s \geq 0, \forall \psi \in \mathcal{H}^{s}, \theta \mapsto U_{\omega}(\theta) \psi \in C^{0}\left(\mathbb{T}^{n} ; \mathcal{H}^{s}\right)$.
(ii) $\forall s \geq 0, \exists C_{s}>0$ such that for all $\theta \in \mathbb{T}^{n}$
\[

$$
\begin{equation*}
\left\|\mathbf{1}-U_{\omega}(\theta)\right\|_{\mathcal{L}\left(\mathcal{H}^{s+2} ; \mathcal{H}^{s}\right)} \leq C_{s} \epsilon . \tag{1.7}
\end{equation*}
$$

\]

(iii) $\forall s, r \geq 0$, the map $\theta \mapsto U_{\omega}(\theta)$ is of class $C^{r}\left(\mathbb{T}^{n} ; \mathcal{L}\left(\mathcal{H}^{s+4 r+2} ; \mathcal{H}^{s}\right)\right)$.

Remark 1.2. Remark that in Theorem 1.1, if the frequencies $\nu_{j}$ are resonant, then the change of coordinates $U_{\omega}$ is close to the identity (in the sense of (1.7)), but the Hamiltonian $H_{\infty}$ is not necessary diagonal. However it is always possible to diagonalize it by means of a metaplectic transformation which is not close to the identity, see Theorem 2.10 and Remark 2.11 below.

Let us denote by $\mathcal{U}_{\epsilon, \omega}(t, \tau)$ the propagator generated by (1.1) such that $\mathcal{U}_{\epsilon, \omega}(\tau, \tau)=\mathbf{1}, \forall \tau \in \mathbb{R}$. An immediate consequence of Theorem 1.1 is that we have a Floquet decomposition:

$$
\begin{equation*}
\mathcal{U}_{\epsilon, \omega}(t, \tau)=U_{\omega}^{*}(\omega t) \mathrm{e}^{-\mathrm{i}(t-\tau) H_{\infty}} U_{\omega}(\omega \tau) \tag{1.8}
\end{equation*}
$$

An other consequence of 1.8 is that for any $s>0$ the norm $\left\|\mathcal{U}_{\epsilon, \omega}(t, 0) \psi_{0}\right\|_{s}$ is bounded uniformly in time:
Corollary 1.3. Let $\omega \in \mathcal{E}_{\epsilon}$ with $|\epsilon|<\epsilon_{*}$. The following is true: for any $s>0$ one has

$$
\begin{equation*}
c_{s}\left\|\psi_{0}\right\|_{s} \leq\left\|\mathcal{U}_{\epsilon, \omega}(t, 0) \psi_{0}\right\|_{s} \leq C_{s}\left\|\psi_{0}\right\|_{s}, \quad \forall t \in \mathbb{R}, \forall \psi_{0} \in \mathcal{H}^{s} \tag{1.9}
\end{equation*}
$$

for some $c_{s}>0, C_{s}>0$.
Moreover there exists a constant $c_{s}^{\prime}$ s.t. if the initial data $\psi_{0} \in \mathcal{H}^{s+2}$ then

$$
\begin{equation*}
\left\|\psi_{0}\right\|_{s}-\epsilon c_{s}^{\prime}\left\|\psi_{0}\right\|_{s+2} \leq\left\|\mathcal{U}_{\epsilon, \omega}(t, 0) \psi_{0}\right\|_{s} \leq\left\|\psi_{0}\right\|_{s}+\epsilon c_{s}^{\prime}\left\|\psi_{0}\right\|_{s+2}, \quad \forall t \in \mathbb{R} \tag{1.10}
\end{equation*}
$$

It is interesting to compare estimate $\sqrt{1.9}$ with the corresponding estimate which can be obtained for more general perturbations $\bar{W}(t, x, D)$. So denote by $\mathcal{U}(t, \tau)$ the propagator of $H_{0}+W(t, x, D)$ with $\mathcal{U}(\tau, \tau)=1$. Then in MR17] it is proved that if $W(t, x, \xi)$ is a real polynomial in $(x, \xi)$ of degree at most 2 , the propagator $\mathcal{U}(t, s)$ exists, belongs to $\mathcal{L}\left(\mathcal{H}^{s}\right) \forall s \geq 0$ and fulfills

$$
\left\|\mathcal{U}(t, 0) \psi_{0}\right\|_{s} \leq e^{C_{s}|t|}\left\|\psi_{0}\right\|_{s}, \quad \forall t \in \mathbb{R}
$$

(the estimate is sharp!). If $W(t, x, \xi)$ is a polynomial of degree at most 1 one has

$$
\left\|\mathcal{U}(t, 0) \psi_{0}\right\|_{s} \leq C_{s}(1+|t|)^{s}\left\|\psi_{0}\right\|_{s}, \quad \forall t \in \mathbb{R}
$$

Thus estimate 1.9 improves dramatically the upper bounds proved in MR17] when the perturbation is small and depends quasiperiodically in time with "good" frequencies.

As a final remark we recall that growth of Sobolev norms can indeed happen if the frequencies $\omega$ are not well chosen. In Appendix A, we show that the Schrödinger equation

$$
\mathrm{i} \dot{\psi}=\left[-\frac{1}{2} \partial_{x x}+\frac{x^{2}}{2}+a x \sin \omega t\right] \psi, \quad x \in \mathbb{R}
$$

(which was already studied by Graffi and Yajima in GY00 who showed that the corresponding Floquet operator has continuous spectrum) exhibits growth of Sobolev norms if and only if $\omega= \pm 1$, which are clearly resonant frequencies. We also slightly generalize the example.

An other example of growth of Sobolev norms for the perturbed harmonic oscillator is given by Delort [Del14]. There the perturbation is a pseudodifferential operator of order 0, periodic in time with resonant frequency $\omega=1$.

Remark 1.4. The uniform time estimate given in (1.9) is similar to the main result obtained in [EK09] for small perturbation of the Laplace operator on the torus $\mathbb{T}^{d}$. Concerning perturbations of harmonic oscillators in $\mathbb{R}^{d}$ most reducibility known results are obtained for $d=1$ excepted in [GP16].

Remark 1.5. In [EK09, GP16] the estimate (1.10] is proved without loss of regularity; this is due to the fact that the perturbations treated in [EK09, GP16] are bounded operators. There are also some cases (see e.g. $\overline{B G 01])}$ in which the reducing transformation is bounded notwithstanding the fact that the perturbation is unbounded, but this is due to the fact that the unperturbed system has suitable gap properties which are not fulfilled in our case.

Remark 1.6. The $\epsilon^{1 / 9}$ estimate on the measure of the set of resonant frequencies is not optimal. We wrote it just for the sake of giving a simple quantitative estimate.

Remark 1.7. Denote by $\left\{\psi_{k}\right\}_{k \in \mathbb{N}^{d}}$ the set of Hermite functions, namely the eigenvectors of $H_{0}: H_{0} \psi_{k}=\lambda_{k} \psi_{k}$. They form an orthonormal basis of $L^{2}\left(\mathbb{R}^{d}\right)$, and writing $\psi=\sum_{k} c_{k} \psi_{k}$ one has $\|\psi\|_{s}^{2} \simeq \sum_{k}(1+|k|)^{2 s}\left|c_{k}\right|^{2}$. Denote $\psi(t)=\sum_{k \in \mathbb{N}^{d}} c_{k}(t) \psi_{k}$ the solution of (1.1) written on the Hermite basis. Then 1.9) implies the following dynamical localization for the energy of the solution: $\forall s \geq 0, \exists C_{s} \equiv C_{s}\left(\psi_{0}\right)>0$ :

$$
\begin{equation*}
\sup _{t \in \mathbb{R}}\left|c_{k}(t)\right| \leq C_{s}(1+|k|)^{-s} \quad, \quad \forall k \in \mathbb{N}^{d} \tag{1.11}
\end{equation*}
$$

From the dynamical property $(1.11)$ one obtains easily that every state $\psi \in L^{2}\left(\mathbb{R}^{d}\right)$ is a bounded state for the time evolution $\mathcal{U}_{\epsilon, \omega}(t, 0) \psi$ under the conditions of Theorem 1.1 on $(\epsilon, \omega)$. The corresponding definitions are given in [EV83]:
Definition 1.8 (See EV83). A function $\psi \in L^{2}\left(\mathbb{R}^{d}\right)$ is a bounded state (or belongs to the point spectral subspace of $\left.\left\{\mathcal{U}_{\epsilon, \omega}(t, 0)\right\}_{t \in \mathbb{R}}\right)$ if the quantum trajectory $\left\{\mathcal{U}_{\epsilon, \omega}(t, 0) \psi: \quad t \in \mathbb{R}\right\}$ is a precompact subset of $L^{2}\left(\mathbb{R}^{d}\right)$.

Corollary 1.9. Under the conditions of Theorem 1.1 on $(\epsilon, \omega)$, every state $\psi \in L^{2}\left(\mathbb{R}^{d}\right)$ is a bounded state of $\left\{\mathcal{U}_{\epsilon, \omega}(t, 0)\right\}_{t \in \mathbb{R}}$.

Proof. To prove that every state $\psi \in L^{2}\left(\mathbb{R}^{d}\right)$ is a bounded state for the time evolution $\mathcal{U}_{\epsilon, \omega}(t, 0) \psi$, using that $\mathcal{H}^{s}$ is dense in $L^{2}\left(\mathbb{R}^{d}\right)$, it is enough to assume that $\psi \in \mathcal{H}^{s}$, with $s>\frac{d}{2}$. With the notations of Remark 1.7, we write

$$
\psi(t)=\psi^{(N)}(t)+R^{(N)}(t)
$$

where $\psi^{(N)}(t)=\sum_{|k| \leq N} c_{k}(t) \psi_{k}$ and $R^{(N)}(t)=\sum_{|k|>N} c_{k}(t) \psi_{k}$.
Take $\delta>0$. Applying (1.7), taking $N$ large enough, we get that for all $t \in \mathbb{R},\left\|R^{(N)}(t)\right\|_{0} \leq \frac{\delta}{2}$. But $\left\{\psi^{(N)}(t), t \in \mathbb{R}\right\}$ is a subset of a finite dimensional linear space. So we get that $\left\{\mathcal{U}_{\epsilon, \omega}(t, 0) \psi\right.$ : $t \in \mathbb{R}\}$ is a precompact subset of $L^{2}\left(\mathbb{R}^{d}\right)$.

This last dynamical result is deeply connected with the spectrum of the Floquet operator. First remark that Theorem 1.1 implies the following
Corollary 1.10. The operator $U_{\omega}$ induces a unitary transformation $L^{2}\left(\mathbb{T}^{n}\right) \otimes L^{2}\left(\mathbb{R}^{d}\right)$ which transforms the Floquet operator $K$, namely

$$
K:=-\mathrm{i} \omega \cdot \frac{\partial}{\partial \theta}+H_{0}+\epsilon W(\theta)
$$

into

$$
-\mathrm{i} \omega \cdot \frac{\partial}{\partial \theta}+H_{\infty}
$$

Thus one has that the spectrum of $K$ is pure point and its eigenvalues are $\lambda_{j}^{\infty}+\omega \cdot k$.
Notice that Enss and Veselic proved that the spectrum of the Floquet operator is pure point if and only if every state is a bounded state EV83, Theorems 2.3 and Theorem 3.2]. So Corollary 1.10 gives another proof of Corollary 1.9 .

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## 2 Proof

To start with we scale the variables $x_{j}$ by defining $x_{j}^{\prime}=\sqrt{\nu_{j}} x_{j}$ so that, defining

$$
h_{j}\left(x_{j}, \xi_{j}\right):=\xi_{j}^{2}+x_{j}^{2}, \quad H_{j}:=-\partial_{x_{j}}^{2}+x_{j}^{2}
$$

one has

$$
\begin{equation*}
h_{0}=\sum_{j=1}^{d} \nu_{j} h_{j}, \quad H_{0}=\sum_{j=1}^{d} \nu_{j} H_{j} \tag{2.1}
\end{equation*}
$$

Remark 2.1. Notice that for any positive definite quadratic Hamiltonian $h$ on $\mathbb{R}^{2 d}$ there exists a symplectic basis such that $h=\sum_{j=1}^{d} \nu_{j} h_{j}$, with $\nu_{j}>0$ for $1 \leq j \leq d$ (see HÖभ).

For convenience in this paper we shall consider the Weyl quantization. The Weyl quantization of a symbol $f$ is the operator $\mathrm{Op}^{w}(f)$, defined as usual as

$$
\mathrm{Op}^{w}(f) u(x)=\frac{1}{(2 \pi)^{d}} \int_{y, \xi \in \mathbb{R}^{d}} e^{\mathrm{i}(x-y) \xi} f\left(\frac{x+y}{2}, \xi\right) u(y) d y d \xi .
$$

Correspondingly we will say that an operator $F=\mathrm{Op}^{w}(f)$ is the Weyl operator with Weyl symbol $f$. Notice that for polynomials $f$ of degree at most 2 in $(x, \xi), \mathrm{Op}^{w}(f)=f(x, D)+$ const, where $D=\mathrm{i}^{-1} \nabla_{x}$.

Most of the times we also use the notation $f^{w}(x, D):=\mathrm{Op}^{w}(f)$. In particular, in equation (1.2) $W\left(\omega t, x,-\mathrm{i} \partial_{x}\right)$ denotes the Weyl operator $W^{w}(\omega t, x, D)$.

Given a Hamiltonian $\chi=\chi(x, \xi)$, we will denote by $\phi_{\chi}^{t}$ the flow of the corresponding classical Hamilton equations.

It is well known that, if $f$ and $g$ are symbols, then the operator $-\mathrm{i}\left[f^{w}(x, D) ; g^{w}(x, D)\right]$ admits a symbol denoted by $\{f ; g\}_{M}$ (Moyal bracket). Two fundamental properties of quadratic polynomial symbols are given by the following well known remarks.

Remark 2.2. If $f$ or $g$ is a polynomial of degree at most 2, then $\{f ; g\}_{M}=\{f ; g\}$, where

$$
\{f ; g\}:=\sum_{j=1}^{d} \frac{\partial f}{\partial x_{j}} \frac{\partial g}{\partial \xi_{j}}-\frac{\partial g}{\partial x_{j}} \frac{\partial f}{\partial \xi_{j}}
$$

is the Poisson Bracket of $f$ and $g$.

Remark 2.3. Let $\chi$ be a polynomial of degree at most 2, then it follows from the previous remark that, for any Weyl operator $f^{w}(x, D)$, the symbol of $e^{\mathrm{i} t \chi^{w}(x, D)} f^{w}(x, D) e^{-\mathrm{i} t \chi^{w}(x, D)}$ is $f \circ \phi_{\chi}^{t}$.
Remark 2.4. If $f$ and $g$ are not quadratic polynomials, then $\{f ; g\}_{M}=\{f ; g\}+$ lower order terms; similar lower order corrections would appear in the symbol of $e^{-\mathrm{i} t \chi^{w}(x, D)} f^{w}(x, D) e^{\mathrm{i} t \chi^{w}(x, D)}$. That is the reason why we restrict here to the case of quadratic perturbations. In order to deal with more general perturbations one needs further ideas which will be developed elsewhere.

Next we need to know how a time dependent transformation transforms a classical and a quantum Hamiltonian. Precisely, consider a 1-parameter family of (Hamiltonian) functions $\chi(t, x, \xi)$ (where $t$ is thought as an external parameter) and denote by $\phi^{\tau}(t, x, \xi)$ the time $\tau$ flow it generates, precisely the solution of

$$
\begin{equation*}
\frac{d x}{d \tau}=\frac{\partial \chi}{\partial \xi}(t, x, \xi), \quad \frac{d \xi}{d \tau}=-\frac{\partial \chi}{\partial x}(t, x, \xi) . \tag{2.2}
\end{equation*}
$$

Consider the time dependent coordinate transformation

$$
\begin{equation*}
(x, \xi)=\phi^{1}\left(t, x^{\prime}, \xi^{\prime}\right):=\left.\phi^{\tau}\left(t, x^{\prime}, \xi^{\prime}\right)\right|_{\tau=1} \tag{2.3}
\end{equation*}
$$

Remark 2.5. Working in the extended phase space in which time and a new momentum conjugated to it are added, it is easy to see that the coordinate transformation (2.3) transforms a Hamiltonian system with Hamiltonian $h$ into a Hamiltonian system with Hamiltonian $h^{\prime}$ given by

$$
\begin{equation*}
h^{\prime}\left(t, x^{\prime}, \xi^{\prime}\right)=h\left(\phi^{1}\left(t, x^{\prime}, \xi^{\prime}\right)\right)-\int_{0}^{1} \frac{\partial \chi}{\partial t}\left(t, \phi^{\tau}\left(t, x^{\prime}, \xi^{\prime}\right)\right) d \tau \tag{2.4}
\end{equation*}
$$

Remark 2.6. If the operator $\chi^{w}(t, x, D)$ is selfadjoint for any fixed $t$, then the transformation

$$
\begin{equation*}
\psi=e^{-\mathrm{i} \chi^{w}(t, x, D)} \psi^{\prime} \tag{2.5}
\end{equation*}
$$

transform $\mathrm{i} \dot{\psi}=H \psi$ into $\mathrm{i} \dot{\psi}^{\prime}=H^{\prime} \psi^{\prime}$ with

$$
\begin{equation*}
H^{\prime}=e^{\mathrm{i} \chi^{w}(t, x, D)} H e^{-\mathrm{i} \chi^{w}(t, x, D)}-\int_{0}^{1} e^{\mathrm{i} \tau \chi^{w}(t, x, D)}\left(\partial_{t} \chi^{w}(t, x, \xi)\right) e^{-\mathrm{i} \tau \chi^{w}(t, x, D)} d \tau \tag{2.6}
\end{equation*}
$$

This is seen by an explicit computation. For example see Lemma 3.2 of Bam17a].
So in view of Remark 2.3 provided that transformation 2.5 is well defined in the quadratic case, the quantum transformed Hamiltonian (2.6) is the exact quantization of the transformed classical Hamiltonian (2.4).

To study the analytic properties of the transformation (2.5 we will use the following simplified version of Theorem 1.2 of MR17] (to which we refer for the proof).
Theorem 2.7 (MR17). Let $H_{0}$ be the Hamiltonian of the harmonic oscillator. If $X$ is an operator symmetric on $\mathcal{H}^{\infty}$ such that $X H_{0}^{-1}$ and $\left[X, H_{0}\right] H_{0}^{-1}$ belong to $\mathcal{L}\left(\mathcal{H}^{s}\right)$ for any $s \geq 0$, then the Schrödinger equation

$$
\mathrm{i} \partial_{\tau} \psi=X \psi
$$

is globally well posed in $\mathcal{H}^{s}$ for any $s$, and its unitary propagator $e^{-\mathrm{i} \tau X}$ belongs to $\mathcal{L}\left(\mathcal{H}^{s}\right), \forall s \geq 0$. Furthermore one has the quantitative estimate

$$
\begin{equation*}
c_{s}\|\psi\|_{s} \leq\left\|e^{-\mathrm{i} \tau X} \psi\right\|_{s} \leq C_{s}\|\psi\|_{s}, \quad \forall \tau \in[0,1] \tag{2.7}
\end{equation*}
$$

where the constants $c_{s}, C_{s}>0$ depend only on $\left\|\left[X, H_{0}^{s}\right] H_{0}^{-s}\right\|_{\mathcal{L}\left(\mathcal{H}^{0}\right)}$.

The properties of the transformation are given by the next lemma and are closely related to the standard properties on the smoothness in time of the semigroup generated by an unbounded operator.

Lemma 2.8. Let $\chi(\rho, x, \xi)$ be a polynomial in $(x, \xi)$ of degree at most 2 with real coefficients depending in a $C^{\infty}$ way on $\rho \in \mathbb{R}^{n}$. Then $\forall \rho \in \mathbb{R}^{n}$, the operator $\chi^{w}(\rho, x, D)$ is selfadjoint in $L^{2}\left(\mathbb{R}^{d}\right)$. Furthermore $\forall s \geq 0, \forall \tau \in \mathbb{R}$ the following holds true:
(i) the map $\rho \mapsto e^{-\mathrm{i} \tau \chi^{w}(\rho, x, D)} \in C^{0}\left(\mathbb{R}^{n}, \mathcal{L}\left(\mathcal{H}^{s+2}, \mathcal{H}^{s}\right)\right)$.
(ii) $\forall \psi \in \mathcal{H}^{s}$, the map $\rho \mapsto e^{-\mathrm{i} \tau \chi^{w}(\rho, x, D)} \psi \in C^{0}\left(\mathbb{R}^{n}, \mathcal{H}^{s}\right)$.
(iii) $\forall r \in \mathbb{N}$ the map $\rho \mapsto e^{-\mathrm{i} \tau \chi^{w}(\rho, x, D)} \in C^{r}\left(\mathbb{R}^{n}, \mathcal{L}\left(\mathcal{H}^{s+4 r+2}, \mathcal{H}^{s}\right)\right)$.
(iv) If the coefficients of $\chi(\rho, x, \xi)$ are uniformly bounded in $\rho \in \mathbb{R}^{n}$ then for any $s>0$ there exist $c_{s}>0, C_{s}>0$ such that we have

$$
c_{s}\|\psi\|_{s} \leq\left\|e^{-\mathrm{i} \tau \chi^{w}(\rho, x, D)} \psi\right\|_{s} \leq C_{s}\|\psi\|_{s}, \quad \forall \rho \in \mathbb{R}^{n}, \forall \tau \in[0,1]
$$

Proof. First we remark that in this lemma the quantity $\rho$ plays the role of a parameter. Since $\chi(\rho, x, \xi)$ is a real valued polynomial in $(x, \xi)$ of degree at most 2 , the operator $\chi^{w}(\rho, x, D)$ is selfadjoint in $L^{2}\left(\mathbb{R}^{d}\right)$, so $\forall \rho \in \mathbb{R}^{n}$ the propagator $e^{-\mathrm{i} \tau \chi^{w}(\rho, x, D)}$ is unitary on $L^{2}\left(\mathbb{R}^{d}\right)$.
To show that $e^{-\mathrm{i} \tau \chi^{w}(\rho, x, D)}$ maps $\mathcal{H}^{s}$ to itself, $\forall s>0, \forall \rho \in \mathbb{R}^{n}$, we apply Theorem 2.7. Indeed since $\chi^{w}(\rho, x, D)$ has a polynomial symbol, then $\chi^{w}(\rho, x, D) H_{0}^{-1}$ and the commutator [ $\left.H_{0}, \chi^{w}(\rho, x, D)\right] H_{0}^{-1}$ belong to $\mathcal{L}\left(\mathcal{H}^{s}\right), \forall s \geq 0$. Item (iv) follows by estimate 2.7) and the remark that $\left\|\left[H_{0}^{s}, \chi^{w}(\rho, x, D)\right] H_{0}^{-s}\right\|_{\mathcal{L}\left(\mathcal{H}^{0}\right)}$ is bounded uniformly in $\rho$.
To prove item (i) we use the Duhamel formula

$$
\begin{equation*}
e^{-\mathrm{i} \tau B}-e^{-\mathrm{i} \tau A}=\mathrm{i} \int_{0}^{\tau} e^{-\mathrm{i}\left(\tau-\tau_{1}\right) A}(A-B) e^{-\mathrm{i} \tau_{1} B} d \tau_{1} \tag{2.8}
\end{equation*}
$$

Then choosing $B=\chi^{w}\left(\rho+\rho^{\prime}, x, D\right), A=\chi^{w}(\rho, x, D)$ one has that $\forall 0 \leq \tau \leq 1$

$$
\left\|e^{-\mathrm{i} \tau \chi^{w}\left(\rho+\rho^{\prime}, x, D\right)}-e^{-\mathrm{i} \tau \chi^{w}(\rho, x, D)}\right\|_{\mathcal{L}\left(\mathcal{H}^{s+2}, \mathcal{H}^{s}\right)} \leq C\left\|\chi^{w}\left(\rho+\rho^{\prime}, x, D\right)-\chi^{w}(\rho, x, D)\right\|_{\mathcal{L}\left(\mathcal{H}^{s+2}, \mathcal{H}^{s}\right)}
$$

This proves item $(i)$. Continuity in item (ii) is deduced by $(i)$ with a standard density argument. Finally item (iii) is proved by induction on $r$ again using the Duhamel formula (2.8).

Remark 2.5, Remark 2.6 and Lemma 2.8 imply the following important proposition.
Proposition 2.9. Let $\chi(t, x, \xi)$ be a polynomial of degree at most 2 in $x$ and $\xi$ with smooth time dependent coefficients. If the transformation (2.3) transforms a classical system with Hamiltonian $h$ into a Hamiltonian system with Hamiltonian $h^{\prime}$, then the transformation (2.5) transforms the quantum system with Hamiltonian $h^{w}$ into the quantum system with Hamiltonian $\left(h^{\prime}\right)^{w}$.

As a consequence, for quadratic Hamiltonians, the quantum KAM theorem will follow from the corresponding classical KAM theorem.

To give the needed result, consider the classical time dependent Hamiltonian

$$
\begin{equation*}
h_{\epsilon}(\omega t, x, \xi):=\sum_{1 \leq j \leq d} \nu_{j} \frac{x_{j}^{2}+\xi_{j}^{2}}{2}+\epsilon W(\omega t, x, \xi) \tag{2.9}
\end{equation*}
$$

with $W$ as in the introduction. The following KAM theorem holds.

Theorem 2.10. Assume that $\nu_{j} \geq \nu_{0}>0$ for $j=1, \cdots, d$ and that $\mathbb{T}^{n} \times \mathbb{R}^{d} \times \mathbb{R}^{d} \ni(\theta, x, \xi) \mapsto$ $W(\theta, x, \xi) \in \mathbb{R}$ is a polynomial in $(x, \xi)$ of degree at most 2 with coefficients which are real analytic functions of $\theta \in \mathbb{T}^{n}$.
Then there exists $\epsilon_{*}>0$ and $C>0$, such that for $|\epsilon|<\epsilon_{*}$ the following holds true:
(i) there exists a closed set $\mathcal{E}_{\epsilon} \subset(0,2 \pi)^{n}$ with $\operatorname{meas}\left((0,2 \pi)^{n} \backslash \mathcal{E}_{\epsilon}\right) \leq C \epsilon^{\frac{1}{9}}$;
(ii) for any $\omega \in \mathcal{E}_{\epsilon}$, there exists an analytic map $\theta \mapsto A_{\omega}(\theta) \in \operatorname{sp}(2 d)$ (symplectic algebrd ${ }^{2}$ of dimension $2 d$ ) and an analytic map $\theta \mapsto V_{\omega}(\theta) \in \mathbb{R}^{2 d}$, such that the change of coordinates

$$
\begin{equation*}
\left(x^{\prime}, \xi^{\prime}\right)=e^{A_{\omega}(\omega t)}(x, \xi)+V_{\omega}(\omega t) \tag{2.10}
\end{equation*}
$$

conjugates the Hamiltonian equations of (2.9) to the Hamiltonian equations of a homogeneous polynomial $h_{\infty}(x, \xi)$ of degree 2 which is positive definite. Finally both $A_{\omega}$ and $V_{\omega}$ are $\epsilon$ close to zero.

Furthermore $h_{\infty}$ can be diagonalized: there exists a matrix $\mathcal{P} \in \operatorname{Sp}(2 d)$ (symplectic group of dimension $2 d$ ) such that, denoting $(y, \eta)=\mathcal{P}(x, \xi)$ we have

$$
\begin{equation*}
h_{\infty} \circ \mathcal{P}^{-1}(y, \eta)=\sum_{j=1}^{d} \nu_{j}^{\infty}\left(y_{j}^{2}+\eta_{j}^{2}\right) \tag{2.11}
\end{equation*}
$$

where $\nu_{j}^{\infty}=\nu_{j}^{\infty}(\omega)$ are defined on $\mathcal{E}_{\epsilon}$ and fulfill the estimates

$$
\begin{equation*}
\left|\nu_{j}^{\infty}-\nu_{j}\right| \leq C \epsilon, \quad j=1, \cdots, d \tag{2.12}
\end{equation*}
$$

Remark 2.11. In general, the matrix $\mathcal{P}$ is not close to identity. However, in case the frequencies $\nu_{j}$ are non resonant, then $\mathcal{P}=\mathbf{1}$.

KAM theory in finite dimensions is nowadays standard. In particular we believe that Theorem 2.10 can be obtained combining the results of Eli88, You99]. However, for the reader convenience and the sake of being self-contained, we add in Section 3 its proof.

Theorem 1.1 follows immediately combining the results of Theorem 2.10 and Proposition 2.9 .
Proof of Theorem 1.1. We see easily that the change of coordinates 2.10 has the form 2.3 with an Hamiltonian $\chi_{\omega}(\omega t, x, \xi)$ which is a polynomial in $(x, \xi)$ of degree at most 2 with real, smooth and uniformly bounded coefficients in $t \in \mathbb{R}$.
Define $U_{\omega}(\omega t)=e^{-\mathrm{i} \chi_{\omega}^{w}(\omega t, x, D)}$. By Proposition 2.9 it conjugates the original equation (1.1) to 1.6) where $H_{\infty}:=\mathrm{Op}^{w}\left(h_{\infty}\right)$.

Furthermore $\theta \mapsto U_{\omega}(\theta)$ fulfills $(i)-(i v)$ of Lemma 2.8 , from which it follows immediately that $\theta \mapsto U_{\omega}(\theta)$ fulfills item $(i)$, (iii) of Theorem 1.1. Concerning item (ii), by Taylor formula the quantity $\left\|\mathbf{1}-U_{\omega}(\theta)\right\|_{\mathcal{L}\left(\mathcal{H}^{s+2}, \mathcal{H}^{s}\right)}$ is controlled by $\left\|\chi_{\omega}^{w}(\theta, x, D)\right\|_{\mathcal{L}\left(\mathcal{H}^{s+2}, \mathcal{H}^{s}\right)}$, from which estimate (1.7) follows.

Finally using the metaplectic representation (see CR12]) and 2.11, there exists a unitary transformation in $L^{2}, \mathcal{R}\left(\mathcal{P}^{-1}\right)$, such that

$$
\mathcal{R}\left(\mathcal{P}^{-1}\right)^{*} H_{\infty} \mathcal{R}\left(\mathcal{P}^{-1}\right)=\sum_{j=1}^{d} \nu_{j}^{\infty}\left(x_{j}^{2}-\partial_{x_{j}}^{2}\right)
$$

[^2]We prove now Corollary 1.3
Proof of Corollary 1.3. Consider first the propagator $e^{-\mathrm{i} t H_{\infty}}$. We claim that

$$
\begin{equation*}
\sup _{t \in \mathbb{R}}\left\|e^{-i t H_{\infty}}\right\|_{\mathcal{L}\left(\mathcal{H}^{s}\right)}<\infty, \quad \forall t \in \mathbb{R} \tag{2.13}
\end{equation*}
$$

Recall that $H_{\infty}=h_{\infty}^{w}(x, D)$ where $h_{\infty}(x, \xi)$ is a positive definite symmetric form which can be diagonalized by a symplectic matrix $\mathcal{P}$. Since $h_{\infty}$ is positive definite, there exist $c_{0}, c_{1}, c_{2}>0$ s.t.

$$
c_{1} h_{0}(x, \xi) \leq c_{0}+h_{\infty}(x, \xi) \leq c_{2}\left(1+h_{0}(x, \xi)\right)
$$

which implies that $C_{1} H_{0} \leq C_{0}+H_{\infty} \leq C_{2}\left(1+H_{0}\right)$ as bilinear form. Thus one has the equivalence of norms

$$
C_{s}^{-1}\|\psi\|_{\mathcal{H}^{s}} \leq\left\|\left(H_{\infty}\right)^{s / 2} \psi\right\|_{L^{2}} \leq C_{s}\|\psi\|_{\mathcal{H}^{s}} .
$$

Then

$$
\left\|e^{-\mathrm{i} t H_{\infty}} \psi_{0}\right\|_{\mathcal{H}^{s}} \leq C_{s}\left\|\left(H_{\infty}\right)^{s / 2} e^{-\mathrm{i} t H_{\infty}} \psi_{0}\right\|_{L^{2}}=C_{s}\left\|\left(H_{\infty}\right)^{s / 2} \psi_{0}\right\|_{L^{2}} \leq C_{s}^{\prime}\left\|\psi_{0}\right\|_{\mathcal{H}^{s}}
$$

which implies (2.13).
Now let $\psi(t)$ be a solution of (1.1). By formula 1.8), $\psi(t)=U_{\omega}^{*}(\omega t) e^{-\mathrm{i} t H_{\infty}} U_{\omega}(0) \psi_{0}$. Then the upper bound in 1.9 follows easily from 2.13 and $\sup _{t}\left\|U_{\omega}(\omega t)\right\|_{\mathcal{L}\left(\mathcal{H}^{s}\right)}<\infty$, which is a consequence of Lemma 2.8 . The lower-bound follows by applying Lemma 2.8 (iv).

Finally estimate 1.10 follows from (1.7).

## 3 A classical KAM result.

In this section we prove Theorem 2.10. We prefer to work in the extended phase space in which we add the angles $\theta \in \mathbb{T}^{n}$ as new variables and their conjugated momenta $I \in \mathbb{R}^{n}$. Furthermore we will use complex variables defined by

$$
z_{j}=\frac{\xi_{j}-\mathrm{i} x_{j}}{\sqrt{2}}
$$

so that our phase space will be $\mathbb{T}^{n} \times \mathbb{R}^{n} \times \mathbb{C}^{d}$, with $\mathbb{C}^{d}$ considered as a real vector space. The symplectic form is $d I \wedge d \theta+\mathrm{i} d z \wedge d \bar{z}$ and the Hamilton equations of a Hamiltonian function $h(\theta, I, z, \bar{z})$ are

$$
\dot{I}=-\frac{\partial h}{\partial \theta}, \quad \dot{\theta}=\frac{\partial h}{\partial I}, \quad \dot{z}=-\mathrm{i} \frac{\partial h}{\partial \bar{z}}
$$

In this framework $h_{0}$ takes the form $h_{0}=\sum_{j=1}^{d} \nu_{j} z_{j} \bar{z}_{j}$ and $W$ takes the form of polynomial in $z, \bar{z}$ of degree two $W(\theta, x, \xi)=q(\theta, z, \bar{z})$. The Hamiltonian system associated with the time dependent Hamiltonian $h_{\epsilon}$ (see 2.9) is then equivalent to the Hamiltonian system associated with the time independent Hamiltonian $\omega \cdot I+h_{\epsilon}$ (written in complex variables) in the extended phase space.

### 3.1 General strategy

Let $h$ be a Hamiltonian in normal form:

$$
\begin{equation*}
h(I, \theta, z, \bar{z})=\omega \cdot I+\langle z, N(\omega) \bar{z}\rangle \tag{3.1}
\end{equation*}
$$

with $N \in \mathcal{M}_{H}$ the set of Hermitian matrix. Notice that at the beginning of the procedure $N$ is diagonal,

$$
N=N_{0}=\operatorname{diag}\left(\nu_{j}, j=1, \cdots, d\right)
$$

and is independent of $\omega$. Let $q \equiv q_{\omega}$ be a polynomial Hamiltonian which takes real values: $q(\theta, z, \bar{z}) \in \mathbb{R}$ for $\theta \in \mathbb{T}^{n}$ and $z \in \mathbb{C}^{d}$. We write

$$
\begin{equation*}
q(\theta, z, \bar{z})=\left\langle z, Q_{z z}(\theta) z\right\rangle+\left\langle z, Q_{z \bar{z}}(\theta) \bar{z}\right\rangle+\left\langle\bar{z}, \bar{Q}_{z z}(\theta) \bar{z}\right\rangle+\left\langle Q_{z}(\theta), z\right\rangle+\left\langle\bar{Q}_{\bar{z}}(\theta), \bar{z}\right\rangle \tag{3.2}
\end{equation*}
$$

where $Q_{z z}(\theta) \equiv Q_{z z}(\omega, \theta)$ and $Q_{z \bar{z}}(\theta) \equiv Q_{z \bar{z}}(\omega, \theta)$ are $d \times d$ complex matrices and $Q_{z}(\theta) \equiv$ $Q_{z}(\theta, \omega)$ is a vector in $\mathbb{C}^{d}$. They all depend analytically on the angle $\theta \in \mathbb{T}_{\sigma}^{n}:=\{x+i y \mid x \in$ $\left.\mathbb{T}^{n}, y \in \mathbb{R}^{n},|y|<\sigma\right\}$. We notice that $Q_{z \bar{z}}$ is Hermitian while $Q_{z z}$ is symmetric. The size of such polynomial function depending analytically on $\theta \in \mathbb{T}_{\sigma}^{n}$ and $C^{1}$ on $\omega \in \mathcal{D}=(0,2 \pi)^{n}$ will be controlled by the norm

$$
[q]_{\sigma}:=\sup _{\substack{|\operatorname{Im} \theta|<\sigma \\ \omega \in \mathcal{D}, j=0,1}}\left\|\partial_{\omega}^{j} Q_{z z}(\omega, \theta)\right\|+\sup _{\substack{|\operatorname{Im} \theta|<\sigma \\ \omega \in \mathcal{D}, j=0,1}}\left\|\partial_{\omega}^{j} Q_{z \bar{z}}(\omega, \theta)\right\|+\sup _{\substack{|\operatorname{Im} \theta|<\sigma \\ \omega \in \mathcal{D}, j=0,1}}\left|\partial_{\omega}^{j} Q_{z}(\omega, \theta)\right|
$$

and we denote by $\mathcal{Q}(\sigma)$ the class of Hamiltonians of the form (3.2) whose norm $[\cdot]_{\sigma}$ is finite.
Let us assume that $[q]_{\sigma}=\mathcal{O}(\epsilon)$. We search for $\chi \equiv \chi_{\omega} \in \mathcal{Q}(\sigma)$ with $[\chi]_{\sigma}=\mathcal{O}(\epsilon)$ such that its time-one flow $\phi_{\chi} \equiv \phi_{\chi}^{t=1}$ (in the extended phase space, of course) transforms the Hamiltonian $h+q$ into

$$
\begin{equation*}
(h+q(\theta)) \circ \phi_{\chi}=h_{+}+q_{+}(\theta), \quad \omega \in \mathcal{D}_{+} \tag{3.3}
\end{equation*}
$$

where $h_{+}=\omega \cdot I+\left\langle z, N_{+} \bar{z}\right\rangle$ is a new normal form, $\epsilon$-close to $h$, the new perturbation $q_{+} \in \mathcal{Q}(\sigma)$ is of $\operatorname{siz} \epsilon^{3} \mathcal{O}\left(\epsilon^{\frac{3}{2}}\right)$ and $\mathcal{D}_{+} \subset \mathcal{D}$ is $\epsilon^{\alpha}$-close to $\mathcal{D}$ for some $\alpha>0$. Notice that all the functions are defined on the whole open set $\mathcal{D}$ but the equalities holds only on $\mathcal{D}_{+}$a subset of $\mathcal{D}$ from which we excised the "resonant parts".
As a consequence of the Hamiltonian structure we have that

$$
(h+q(\theta)) \circ \phi_{\chi}=h+\{h, \chi\}+q(\theta)+\mathcal{O}\left(\epsilon^{\frac{3}{2}}\right), \quad \omega \in \mathcal{D}_{+} .
$$

So to achieve the goal above we should solve the homological equation:

$$
\begin{equation*}
\{h, \chi\}=h_{+}-h-q(\theta)+\mathcal{O}\left(\epsilon^{\frac{3}{2}}\right), \quad \omega \in \mathcal{D}_{+} \tag{3.4}
\end{equation*}
$$

Repeating iteratively the same procedure with $h_{+}$instead of $h$, we will construct a change of variable $\phi$ such that

$$
(h+q(\theta)) \circ \phi=\omega \cdot I+h_{\infty}, \quad \omega \in \mathcal{D}_{\infty},
$$

with $h_{\infty}=\left\langle z, N_{\infty}(\omega) \bar{z}\right\rangle$ in normal form and $\mathcal{D}_{\infty}$ a $\epsilon^{\alpha}$-close subset of $\mathcal{D}$. Note that we will be forced to solve the homological equation not only for the diagonal normal form $N_{0}$, but for more general normal form Hamiltonians (3.1) with $N$ close to $N_{0}$.

[^3]
### 3.2 Homological equation

Proposition 3.1. Let $\mathcal{D}=(0,2 \pi)^{n}$ and $\mathcal{D} \ni \omega \mapsto N(\omega) \in \mathcal{M}_{H}$ be a $C^{1}$ mapping that verifies

$$
\begin{equation*}
\left\|\partial_{\omega}^{j}\left(N(\omega)-N_{0}\right)\right\|<\frac{\min \left(1, \nu_{0}\right)}{\max (4, d)} \tag{3.5}
\end{equation*}
$$

for $j=0,1$ and $\omega \in \mathcal{D}$. Let $h=\omega \cdot I+\langle z, N \bar{z}\rangle, q \in \mathcal{Q}(\sigma), \kappa>0$ and $K \geq 1$.
Then there exists a closed subset $\mathcal{D}^{\prime}=\mathcal{D}^{\prime}(\kappa, K) \subset \mathcal{D}$, satisfying

$$
\begin{equation*}
\operatorname{meas}\left(\mathcal{D} \backslash \mathcal{D}^{\prime}\right) \leq C K^{n} \kappa, \tag{3.6}
\end{equation*}
$$

and there exist $\chi, r \in \cap_{0 \leq \sigma^{\prime}<\sigma} \mathcal{Q}\left(\sigma^{\prime}\right)$ and $\mathcal{D} \ni \omega \mapsto \tilde{N}(\omega) \in \mathcal{M}_{H}$ a $C^{1}$ mapping such that for all $\omega \in \mathcal{D}^{\prime}$

$$
\begin{equation*}
\{h, \chi\}+q=\langle z, \tilde{N} \bar{z}\rangle+r . \tag{3.7}
\end{equation*}
$$

Furthermore for all $\omega \in \mathcal{D}$

$$
\begin{equation*}
\left\|\partial_{\omega}^{j} \tilde{N}(\omega)\right\| \leq[q]_{\sigma}, \quad j=0,1 \tag{3.8}
\end{equation*}
$$

and for all $0 \leq \sigma^{\prime}<\sigma$

$$
\begin{align*}
{[r]_{\sigma^{\prime}} } & \leq C \frac{e^{-\frac{1}{2}\left(\sigma-\sigma^{\prime}\right) K}}{\left(\sigma-\sigma^{\prime}\right)^{n}}[q]_{\sigma},  \tag{3.9}\\
{[\chi]_{\sigma^{\prime}} } & \leq \frac{C K}{\kappa^{2}\left(\sigma-\sigma^{\prime}\right)^{n}}[q]_{\sigma} \tag{3.10}
\end{align*}
$$

Proof. Writing the Hamiltonians $h, q$ and $\chi$ as in (3.2), the homological equation (3.7) is equivalent to the three following equations (we use that $N$ is Hermitian, thus $\bar{N}={ }^{t} N$ ):

$$
\begin{gather*}
\omega \cdot \nabla_{\theta} X_{z \bar{z}}-\mathrm{i}\left[N, X_{z \bar{z}}\right]=\tilde{N}-Q_{z \bar{z}}+R_{z \bar{z}}  \tag{3.11}\\
\omega \cdot \nabla_{\theta} X_{z z}-\mathrm{i}\left(N X_{z z}+X_{z z} \bar{N}\right)=-Q_{z z}+R_{z z}  \tag{3.12}\\
\omega \cdot \nabla_{\theta} X_{z}+\mathrm{i} N X_{z}=-Q_{z}+R_{z} \tag{3.13}
\end{gather*}
$$

First we solve (3.11). To simplify notations we drop the indices $z \bar{z}$. Written in Fourier variables (w.r.t. $\theta$ ), (3.11) reads

$$
\begin{equation*}
\mathrm{i} \omega \cdot k \hat{X}_{k}-\mathrm{i}\left[N, \hat{X}_{k}\right]=\delta_{k, 0} \tilde{N}-\hat{Q}_{k}+\hat{R}_{k}, \quad k \in \mathbb{Z}^{n} \tag{3.14}
\end{equation*}
$$

where $\delta_{k, j}$ denotes the Kronecker symbol.
When $k=0$ we solve this equation by defining

$$
\hat{X}_{0}=0, \quad \hat{R}_{0}=0 \quad \text { and } \quad \tilde{N}=\hat{Q}_{0} .
$$

We notice that $\tilde{N} \in \mathcal{M}_{H}$ and satisfies (3.8).
When $|k| \geq K$ equation (3.14) is solved by defining

$$
\begin{equation*}
\hat{R}_{k}=\hat{Q}_{k}, \quad \hat{X}_{k}=0 \text { for }|k| \geq K \tag{3.15}
\end{equation*}
$$

Then we set

$$
\hat{R}_{k}=0 \quad \text { for } \quad|k| \leq K
$$

in such a way that $r \in \cap_{0 \leq \sigma^{\prime}<\sigma \mathcal{Q}}\left(\sigma^{\prime}\right)$ and by a standard argument $r$ satisfies 3.9. Now it remains to solve the equations for $\hat{X}_{k}, 0<|k| \leq K$ which we rewrite as

$$
\begin{equation*}
L_{k}(\omega) \hat{X}_{k}=i \hat{Q}_{k} \tag{3.16}
\end{equation*}
$$

where $L_{k}(\omega)$ is the linear operator from $\mathcal{M}_{S}$, the space of symmetric matrices, into itself defined by

$$
L_{k}(\omega): M \mapsto(k \cdot \omega) M-[N(\omega), M] .
$$

We notice that $\mathcal{M}_{S}$ can be endowed with the Hermitian product: $(A, B)=\operatorname{Tr}(\bar{A} B)$ associated with the Hilbert Schmidt norm. Since $N$ is Hermitian, $L_{k}(\omega)$ is self adjoint for this structure. As a first consequence we get

$$
\begin{equation*}
\left\|\left(L_{k}(\omega)\right)^{-1}\right\| \leq \frac{1}{\min \left\{|\lambda|, \lambda \in \Sigma\left(L_{k}(\omega)\right)\right\}}=\frac{1}{\min \{|k \cdot \omega-\alpha(\omega)+\beta(\omega)| \mid \alpha, \beta \in \Sigma(N(\omega))\}} \tag{3.17}
\end{equation*}
$$

where for any matrix $A$, we denote its spectrum by $\Sigma(A)$.
Let us recall an important result of perturbation theory which is a consequence of Theorem 1.10 in Kat95 (since hermitian matrices are normal matrices):

Theorem 3.2 (Kat95] Theorem 1.10). Let $I \subset \mathbb{R}$ and $I \ni z \mapsto M(z)$ a holomorphic curve of hermitian matrices. Then all the eigenvalues and associated eigenvectors of $M(z)$ can be parametrized holomorphically on I.

Let us assume for a while that $N$ depends analytically of $\omega$ in such a way that $\omega \mapsto L_{k}(\omega)$ is analytic. Fix a direction $z_{k} \in \mathbb{R}^{n}$, the eigenvalue $\lambda_{k}(\omega)=k \cdot \omega-\alpha(\omega)+\beta(\omega)$ of $L_{k}(\omega)$ is $C^{1}$ in the directior ${ }^{4} z_{k}$ and the associated unitary eigenvector, denoted by $v(\omega)$, is also piece wise $C^{1}$ in the direction $z_{k}$. Then, as a consequence of the hermiticity of $L_{k}(\omega)$ we have

$$
\partial_{\omega} \lambda(\omega) \cdot z_{k}=\left\langle v(\omega),\left(\partial_{\omega} L_{k}(\omega) \cdot z_{k}\right) v(\omega)\right\rangle .
$$

Therefore, if $N$ depends analytically of $\omega$, we deduce using 3.5 and choosing $z_{k}=\frac{k}{|k|}$

$$
\begin{equation*}
\left|\partial_{\omega} \lambda_{k}(\omega) \cdot \frac{k}{|k|}\right| \geq|k|-2\left\|\partial_{\omega} N\right\| \geq \frac{1}{2} \quad \text { for } k \neq 0 \tag{3.18}
\end{equation*}
$$

which extends also to the points of discontinuity of $v(\omega)$. Now given a matrix $L$ depending on the parameter $\omega \in \mathcal{D}$, we define

$$
\mathcal{D}(L, \kappa)=\left\{\omega \in \mathcal{D} \mid\left\|L(\omega)^{-1}\right\| \leq \kappa^{-1}\right\}
$$

and we recall the following classical lemma:
Lemma 3.3. Let $f:[0,1] \mapsto \mathbb{R}$ a $C^{1}$-map satisfying $\left|f^{\prime}(x)\right| \geq \delta$ for all $x \in[0,1]$ and let $\kappa>0$. Then

$$
\operatorname{meas}\left\{x \in[0,1]||f(x)| \leq \kappa\} \leq \frac{\kappa}{\delta}\right.
$$

Combining this Lemma 3.17 and 3.18 we deduce that, if $N$ depends analytically of $\omega$, then for $k \neq 0$

$$
\begin{equation*}
\operatorname{meas}\left(\mathcal{D} \backslash \mathcal{D}\left(L_{k}, \kappa\right)\right) \leq C \kappa \tag{3.19}
\end{equation*}
$$

[^4]Now it turns out that, by a density argument, this last estimate remains valid (with a larger constant $C$ ) when $N$ is only a $C^{1}$ function of $\omega$ : the point is that (3.18) holds true uniformly for close analytic approximations of $N$.
In particular defining

$$
\mathcal{D}^{\prime}=\bigcap_{0<|k| \leq K} \mathcal{D}\left(L_{k}, \kappa\right)
$$

$\mathcal{D}^{\prime}$ is closed and satisfies (3.6).
By construction, $\hat{X}_{k}(\omega):=i L_{k}(\omega)^{-1} \hat{Q}_{k}$ satisfies 3.16 for $0<|k| \leq K$ and $\omega \in \mathcal{D}\left(L_{k}, \kappa\right)$ and

$$
\begin{equation*}
\left\|\hat{X}_{k}(\omega)\right\| \leq \kappa^{-1}\left\|\hat{Q}_{k}(\omega)\right\|, \quad \omega \in \mathcal{D}\left(L_{k}, \kappa\right) \tag{3.20}
\end{equation*}
$$

It remains to extend $\hat{X}_{k}(\cdot)$ on $\mathcal{D}$. Using again (3.5) we have for any $|k| \leq K$ and any unit vector $z,\left|\partial_{\omega} \lambda(\omega) \cdot z\right| \leq C K$. Therefore

$$
\operatorname{dist}\left(\mathcal{D} \backslash \mathcal{D}\left(L_{k}, \kappa\right), \mathcal{D}\left(L_{k}, \kappa / 2\right)\right) \geq \frac{\kappa}{C K}
$$

and we can construct (by a convolution argument) for each $k, 0<|k| \leq K$, a $C^{1}$ function $g_{k}$ on $\mathcal{D}$ with

$$
\begin{equation*}
\left|g_{k}\right|_{C^{0}(\mathcal{D})} \leq C, \quad\left|g_{k}\right|_{C^{1}(\mathcal{D})} \leq C K \kappa^{-1} \tag{3.21}
\end{equation*}
$$

(the constant C is independent of k ) and such that $g_{k}(\omega)=1$ for $\omega \notin \mathcal{D}\left(L_{k}, \kappa\right)$ and $g_{k}(\omega)=0$ for $\omega \in \mathcal{D}\left(L_{k}, \kappa / 2\right)$. Then $\tilde{X}_{k}=g_{k} \hat{X}_{k}$ is a $C^{1}$ extension of $\hat{X}_{k}$ to $\mathcal{D}$. Similarly we define $\tilde{Q}_{k},=g_{k} \hat{Q}_{k}$ in such a way that $\tilde{X}_{k}$ satisfies

$$
L_{k}(\omega) \tilde{X}_{k}(\omega)=i \tilde{Q}_{k}(\omega), \quad 0<|k| \leq K, \omega \in \mathcal{D}
$$

Differentiating with respect to $\omega$ leads to

$$
L_{k}(\omega) \partial_{\omega_{j}} \hat{X}(k)=\mathrm{i} \partial_{\omega_{j}} \hat{Q}(k)-k_{j} \hat{X}(k)+\left[\partial_{\omega_{j}} N, \hat{X}(k)\right], \quad 1 \leq j \leq n
$$

Denoting $B_{k}(\omega)=\mathrm{i} \partial_{\omega_{j}} \tilde{Q}_{k}(\omega)-k_{j} \tilde{X}_{k}(\omega)+\left[\partial_{\omega_{j}} N(\omega), \tilde{X}_{k}(\omega)\right]$ we have

$$
\left\|\partial_{\omega_{j}} \tilde{X}_{k}(\omega)\right\| \leq \kappa^{-1}\left\|B_{k}(\omega)\right\|, \quad \omega \in \mathcal{D}
$$

Using (3.5, 3.20) and (3.21) we get for $|k| \leq K$ and $\omega \in \mathcal{D}$

$$
\begin{aligned}
\left\|B_{k}(\omega)\right\| & \leq\left\|\partial_{\omega_{j}} \tilde{Q}_{k}(\omega)\right\|+K\left\|\tilde{X}_{k}(\omega)\right\|+2\left\|\partial_{\omega_{j}} N(\omega)\right\|\left\|\tilde{X}_{k}(\omega)\right\| \\
& \leq C K \kappa^{-1}\left(\left\|\partial_{\omega_{j}} \hat{Q}(k, \omega)\right\|+\|\hat{Q}(k, \omega)\|\right) .
\end{aligned}
$$

Combining the last two estimates we get

$$
\sup _{\omega \in \mathcal{D}, j=0,1}\left\|\partial_{\omega}^{j} \tilde{X}_{k}(\omega)\right\| \leq C K \kappa^{-2} \sup _{\omega \in \mathcal{D}, j=0,1}\left\|\partial_{\omega}^{j} \hat{Q}_{k}(\omega)\right\|
$$

Thus defining

$$
X_{z \bar{z}}(\omega, \theta)=\sum_{0<|k| \leq K} \tilde{X}_{k}(\omega) e^{i k \cdot \theta}
$$

$X_{z \bar{z}}(\omega, \cdot)$ satisfies (3.11) for $\omega \in \mathcal{D}^{\prime}$ and leads to 3.10) for $\chi_{z \bar{z}}(\omega, \theta, z, \bar{z})=\left\langle z, X_{z \bar{z}}(\omega, \cdot) \bar{z}\right\rangle$.
We solve 3.13 in a similar way. We notice that in this case we face the small divisors $|\omega \cdot k-\alpha(\omega)|, k \in \mathbb{Z}^{n}$ where $\alpha \in \Sigma(N(\omega))$. In particular for $k=0$ these quantities are $\geq \frac{\nu_{0}}{2}$ since $\left|\alpha-\nu_{j}\right| \leq \frac{\nu_{0}}{4}$ for some $1 \leq j \leq d$ by 3.5.

Written in Fourier and dropping indices $z z 3$ reads

$$
\begin{equation*}
\mathrm{i} \omega \cdot k \hat{X}(k)-\mathrm{i}(N \hat{X}(k)+\hat{X}(k) \bar{N})=-\hat{Q}(k)+\hat{R}(k) . \tag{3.22}
\end{equation*}
$$

So to mimic the resolution of (3.14) we have to replace the operator $L_{k}(\omega)$ by the operator $M_{k}(\omega)$ defined on $\mathcal{M}_{S}$ by

$$
M_{k}(\omega) X:=\omega \cdot k+N X+X \bar{N}
$$

This operator is still self adjoint for the Hermitian product $(A, B)=\operatorname{Tr}(\bar{A} B)$ so the same strategy apply. Nevertheless we have to consider differently the case $k=0$. In that case we use that the eigenvalues of $M_{0}(\omega)$ are close to eigenvalues of the operator $M_{0}$ defined by

$$
M_{0}: X \mapsto N_{0} X+X \bar{N}_{0}=N_{0} X+X N_{0}
$$

with $N_{0}=\operatorname{diag}\left(\nu_{j}, j=1, \cdots, d\right)$ a real and diagonal matrix. Actually in view of 3.5

$$
\left\|\left(L-L_{0}\right) M\right\|_{H S} \leq\left\|N-N_{0}\right\|_{H S}\|M\|_{H S} \leq d\left\|N-N_{0}\right\|\|M\|_{H S} \leq \nu_{0}
$$

The eigenvalues of $L_{0}$ are $\left\{\nu_{j}+\nu_{\ell} \mid j, \ell=1, \cdots, d\right\}$ and they are all larger than $2 \nu_{0}$. We conclude that all the eigenvalues of $M_{0}(\omega)$ satisfy $|\alpha(\omega)| \geq \nu_{0}$. The end of the proof follow as before.

### 3.3 The KAM step.

Theorem 2.10 is proved by an iterative KAM procedure. We begin with the initial Hamiltonian $h_{0}+q_{0}$ where

$$
\begin{equation*}
h_{0}(I, \theta, z, \bar{z})=\omega \cdot I+\left\langle z, N_{0} \bar{z}\right\rangle \tag{3.23}
\end{equation*}
$$

$N_{0}=\operatorname{diag}\left(\nu_{j}, j=1, \cdots, d\right), \omega \in \mathcal{D} \equiv[1,2]^{n}$ and the quadratic perturbation $q_{0}=\epsilon W \in \mathcal{Q}(\sigma, \mathcal{D})$ for some $\sigma>0$. Then we construct iteratively the change of variables $\phi_{m}$, the normal form $h_{m}=\omega \cdot I+\left\langle z, N_{m} \bar{z}\right\rangle$ and the perturbation $q_{m} \in \mathcal{Q}\left(\sigma_{m}, \mathcal{D}_{m}\right)$ as follows: assume that the construction is done up to step $m \geq 0$ then
(i) Using Proposition 3.1 we construct $\chi_{m+1}, r_{m+1}$ and $\tilde{N}_{m}$ the solution of the homological equation:

$$
\begin{equation*}
\left\{h, \chi_{m+1}\right\}=\left\langle z, \tilde{N}_{m} \bar{z}\right\rangle-q_{m}(\theta)+r_{m+1}, \quad \omega \in \mathcal{D}_{m+1}, \theta \in \mathbb{T}_{\sigma_{m+1}}^{n} \tag{3.24}
\end{equation*}
$$

(ii) We define $h_{m+1}:=\omega \cdot I+\left\langle z, N_{m+1} \bar{z}\right\rangle$ by

$$
\begin{equation*}
N_{m+1}=N_{m}+\tilde{N}_{m} \tag{3.25}
\end{equation*}
$$

and

$$
\begin{equation*}
q_{m+1}:=r_{m}+\int_{0}^{1}\left\{(1-t)\left(h_{m+1}-h_{m}+r_{m+1}\right)+t q_{m}, \chi_{m+1}\right\} \circ \phi_{\chi_{m+1}}^{t} \mathrm{~d} t \tag{3.26}
\end{equation*}
$$

By construction, if $Q_{m}$ and $N_{m}$ are Hermitian, so are $R_{m}$ and $S_{m+1}$ by the resolution of the homological equation, and also $N_{m+1}$ and $Q_{m+1}$.
For any regular Hamiltonian $f$ we have, using the Taylor expansion of $f \circ \phi_{\chi_{m+1}}^{t}$ between $t=0$ and $t=1$

$$
f \circ \phi_{\chi_{m+1}}^{1}=f+\left\{f, \chi_{m+1}\right\}+\int_{0}^{1}(1-t)\left\{\left\{f, \chi_{m+1}\right\}, \chi_{m+1}\right\} \circ \phi_{\chi_{m+1}}^{t} \mathrm{~d} t
$$

Therefore we get for $\omega \in \mathcal{D}_{m+1}$

$$
\left(h_{m}+q_{m}\right) \circ \phi_{\chi_{m+1}}^{1}=h_{m+1}+q_{m+1}
$$

### 3.4 Iterative lemma

Following the general scheme above we have

$$
\left(h_{0}+q_{0}\right) \circ \phi_{\chi_{1}}^{1} \circ \cdots \circ \phi_{\chi_{m}}^{1}=h_{m}+q_{m}
$$

where $q_{m}$ is a polynomial of degree two and $h_{m}=\omega \cdot I+\left\langle z, N_{m} \bar{z}\right\rangle$ with $N_{m}$ a Hermitian matrix. At step $m$ the Fourier series are truncated at order $K_{m}$ and the small divisors are controlled by $\kappa_{m}$. Now we specify the choice of all the parameters for $m \geq 0$ in term of $\epsilon_{m}$ which will control $\left[q_{m}\right]_{\mathcal{D}_{m}, \sigma_{m}}$.
First we define $\epsilon_{0}=\epsilon, \sigma_{0}=\sigma, \mathcal{D}_{0}=\mathcal{D}$ and for $m \geq 1$ we choose

$$
\sigma_{m-1}-\sigma_{m}=C_{*} \sigma_{0} m^{-2}, \quad K_{m}=2\left(\sigma_{m-1}-\sigma_{m}\right)^{-1} \ln \epsilon_{m-1}^{-1}, \quad \kappa_{m}=\epsilon_{m-1}^{\frac{1}{8}}
$$

where $\left(C_{*}\right)^{-1}=2 \sum_{j \geq 1} \frac{1}{j^{2}}$.

Lemma 3.4. There exists $\epsilon_{*}>0$ depending on $d$, $n$ such that, for $|\epsilon| \leq \epsilon_{*}$ and

$$
\epsilon_{m}=\epsilon^{(3 / 2)^{m}} \quad, \quad m \geq 0
$$

we have the following:
For all $m \geq 1$ there exist closed subsets $\mathcal{D}_{m} \subset \mathcal{D}_{m-1}, h_{m}=\omega \cdot I+\left\langle z, N_{m} \bar{z}\right\rangle$ in normal form where $\mathcal{D}_{m} \ni \omega \mapsto N_{m}(\omega) \in \mathcal{M}_{H} \in C^{1}$ and there exist $\chi_{m}, q_{m} \in \mathcal{Q}\left(\mathcal{D}_{m}, \sigma_{m}\right)$ such that for $m \geq 1$
(i) The symplectomorphism

$$
\begin{equation*}
\phi_{m} \equiv \phi_{\chi_{m}}(\omega): \mathbb{R}^{n} \times \mathbb{T}^{n} \times \mathbb{C}^{2 d} \rightarrow \mathbb{R}^{n} \times \mathbb{T}^{n} \times \mathbb{C}^{2 d}, \quad \omega \in \mathcal{D}_{m} \tag{3.27}
\end{equation*}
$$

is an affine transformation in $(z, \bar{z})$, analytic in $\theta \in \mathbb{T}_{\sigma_{m}}^{n}$ and $C^{1}$ in $\omega \in \mathcal{D}_{m}$ of the form

$$
\begin{equation*}
\phi_{m}(I, \theta, z, \bar{z})=\left(g_{m}(I, \theta, z, \bar{z}), \theta, \Psi_{m}(\theta, z, \bar{z})\right), \tag{3.28}
\end{equation*}
$$

where, for each $\theta \in \mathbb{T}^{n},(z, \bar{z}) \mapsto \Psi_{m}(\theta, z, \bar{z})$ is a symplectic change of variable on $\mathbb{C}^{2 n}$. The map $\phi_{m}$ links the Hamiltonian at step $m-1$ and the Hamiltonian at step m, i.e.

$$
\left(h_{m-1}+q_{m-1}\right) \circ \phi_{m}=h_{m}+q_{m}, \quad \forall \omega \in \mathcal{D}_{m}
$$

(ii) We have the estimates

$$
\begin{align*}
& \operatorname{meas}\left(\mathcal{D}_{m-1} \backslash \mathcal{D}_{m}\right) \leq \epsilon_{m-1}^{\frac{1}{9}}  \tag{3.29}\\
& {\left[\tilde{N}_{m-1}\right]_{s, \beta}^{\mathcal{D}_{m}} } \leq \epsilon_{m-1},  \tag{3.30}\\
& {\left[q_{m}\right]_{s, \beta}^{\mathcal{D}_{m}, \sigma_{m}} } \leq \epsilon_{m},  \tag{3.31}\\
&\left\|\phi_{m}(\omega)-\mathbf{1}\right\|_{\mathcal{L}\left(\mathbb{R}^{n} \times \mathbb{T}^{n} \times \mathbb{C}^{2 d}\right)} \leq C \epsilon_{m-1}^{\frac{1}{2}}, \quad \forall \omega \in \mathcal{D}_{m} \tag{3.32}
\end{align*}
$$

Proof. At step $1, h_{0}=\omega \cdot I+\left\langle z, N_{0} \bar{z}\right\rangle$ and thus hypothesis (3.5) is trivially satisfied and we can apply Proposition 3.1 to construct $\chi_{1}, N_{1}, r_{1}$ and $\mathcal{D}_{1}$ such that for $\omega \in \mathcal{D}_{1}$

$$
\left\{h_{0}, \chi_{1}\right\}=\left\langle z,\left(N_{1}-N_{0}\right) \bar{z}\right\rangle-q_{0}+r_{1}
$$

Then, using 3.6, we have

$$
\operatorname{meas}\left(\mathcal{D} \backslash \mathcal{D}_{1}\right) \leq C K_{1}^{n} \kappa_{1} \leq \epsilon_{0}^{\frac{1}{9}}
$$

for $\epsilon=\epsilon_{0}$ small enough. Using 3.10 we have for $\epsilon_{0}$ small enough

$$
\left[\chi_{1}\right]_{\mathcal{D}_{1}, \sigma_{1}} \leq C \frac{K_{1}}{\kappa_{1}^{2}\left(\sigma_{0}-\sigma_{1}\right)^{n}} \epsilon_{0} \leq \epsilon_{0}^{\frac{1}{2}}
$$

Similarly using (3.9), (3.8) we have

$$
\left\|N_{1}-N_{0}\right\| \leq \epsilon_{0}
$$

and

$$
\left[r_{1}\right]_{\mathcal{D}_{1}, \sigma_{1}} \leq C \frac{\epsilon_{0}^{\frac{15}{8}}}{\left(\sigma_{1}-\sigma_{0}\right)^{n}} \leq \epsilon_{0}^{\frac{7}{4}}
$$

for $\epsilon=\epsilon_{0}$ small enough. In particular we deduce $\left\|\phi_{1}-\mathbf{1}\right\|_{\mathcal{L}\left(\mathbb{R}^{n} \times \mathbb{T}^{n} \times \mathbb{C}^{2 d}\right)} \leq \epsilon_{0}^{\frac{1}{2}}$. Thus using (3.26) we get for $\epsilon_{0}$ small enough

$$
\left[q_{1}\right]_{\mathcal{D}_{1}, \sigma_{1}} \leq \epsilon_{0}^{3 / 2}=\epsilon_{1}
$$

The form of the flow (3.28) follows since $\chi_{1}$ is a Hamiltonian of the form (3.2).
Now assume that we have verified Lemma 3.4 up to step $m$. We want to perform the step $m+1$. We have $h_{m}=\omega \cdot I+\left\langle z, N_{m} \bar{z}\right\rangle$ and since

$$
\left\|N_{m}-N_{0}\right\| \leq\left\|N_{m}-N_{0}\right\|+\cdots+\left\|N_{1}-N_{0}\right\| \leq \sum_{j=0}^{m-1} \epsilon_{j} \leq 2 \epsilon_{0}
$$

hypothesis 3.5 is satisfied and we can apply Proposition 3.1 to construct $\mathcal{D}_{m+1}, \chi_{m+1}$ and $q_{m+1}$. Estimates $3.29-3.32$ at step $m+1$ are proved as we have proved the corresponding estimates at step 1.

### 3.5 Transition to the limit and proof of Theorem 2.10

Let $\mathcal{E}_{\epsilon}=\cap_{m \geq 0} \mathcal{D}_{m}$. In view of 3.29, this is a closed set satisfying

$$
\operatorname{meas}\left(\mathcal{D} \backslash \mathcal{E}_{\epsilon}\right) \leq \sum_{m \geq 0} \epsilon_{m}^{\frac{1}{9}} \leq 2 \epsilon_{0}^{\frac{1}{9}}
$$

Let us denote $\widetilde{\phi}_{N}=\phi_{1} \circ \cdots \circ \phi_{N}$. Due to 3.32 it satisfies for $M \leq N$ and for $\omega \in \mathcal{E}-\epsilon$

$$
\left\|\widetilde{\phi}_{N}-\widetilde{\phi}_{M}\right\|_{\mathcal{L}\left(\mathbb{R}^{n} \times \mathbb{T}^{n} \times \mathbb{C}^{2 d}\right)} \leq \sum_{m=M}^{N} \epsilon_{m}^{\frac{1}{2}} \leq 2 \epsilon_{M}^{\frac{1}{2}}
$$

Therefore $\left(\widetilde{\phi}_{N}\right)_{N}$ is a Cauchy sequence in $\mathcal{L}\left(\mathbb{R}^{n} \times \mathbb{T}^{n} \times \mathbb{C}^{2 d}\right)$. Thus when $N \rightarrow \infty$ the mappings $\widetilde{\phi}_{N}$ converge to a limit mapping $\phi_{\infty} \in \mathcal{L}\left(\mathbb{R}^{n} \times \mathbb{T}^{n} \times \mathbb{C}^{2 d}\right)$. Furthermore since the convergence is uniform on $\omega \in \mathcal{E}_{\epsilon}$ and $\theta \in \mathbb{T}_{\sigma / 2}, \phi_{\infty}^{1}$ depends analytically on $\theta$ and $C^{1}$ in $\omega$. Moreover,

$$
\begin{equation*}
\left\|\phi_{\infty}-\mathbf{1}\right\|_{\mathcal{L}\left(\mathbb{R}^{n} \times \mathbb{T}^{n} \times \mathbb{C}^{2 d}\right)} \leq \epsilon_{0}^{\frac{1}{2}} \tag{3.33}
\end{equation*}
$$

By construction, the map $\widetilde{\phi}_{m}$ transforms the original Hamiltonian $h_{0}+q_{0}$ into $h_{m}+q_{m}$. When $m \rightarrow \infty$, by 3.31 we get $q_{m} \rightarrow 0$ and by 3.30 we get $N_{m} \rightarrow N$ where

$$
\begin{equation*}
N \equiv N(\omega)=N_{0}+\sum_{k=1}^{+\infty} \tilde{N}_{k} \tag{3.34}
\end{equation*}
$$

is a Hermitian matrix which is $C^{1}$ with respect to $\omega \in \mathcal{E}_{\epsilon}$. Denoting $h_{\infty}(z, \bar{z})=\omega \cdot I+\langle z, N(\omega) \bar{z}\rangle$ we have proved

$$
\begin{equation*}
(h+q(\theta)) \circ \phi_{\infty}=h_{\infty} . \tag{3.35}
\end{equation*}
$$

Furthermore $\forall \omega \in \mathcal{E}_{\epsilon}$ we have, using 3.30,

$$
\left\|N(\omega)-N_{0}\right\| \leq \sum_{m=0}^{\infty} \epsilon_{m} \leq 2 \epsilon
$$

and thus the eigenvalues of $N(\omega)$, denoted $\nu_{j}^{\infty}(\omega)$ satisfy 2.12.
It remains to explicit the affine symplectomorphism $\phi_{\infty}$. At each step of the KAM procedure we have by Lemma 3.4

$$
\phi_{m}(I, \theta, z, \bar{z})=\left(g_{m}(I, \theta, z, \bar{z}), \theta, \Psi_{m}(\theta, z, \bar{z})\right)
$$

and therefore

$$
\phi_{\infty}(I, \theta, z, \bar{z})=(g(I, \theta, z, \bar{z}), \theta, \Psi(\theta, z, \bar{z}))
$$

where $\Psi(\theta, z, \bar{z})=\lim _{m \rightarrow \infty} \Psi_{1} \circ \Psi_{2} \circ \cdots \circ \Psi_{m}$.
It is useful to go back to real variables $(x, \xi)$. More precisely write each Hamiltonian $\chi_{m}$ constructed in the KAM iteration in the variables $(x, \xi)$ :

$$
\chi_{m}(\theta, x, \xi)=\frac{1}{2}\binom{x}{\xi} \cdot E B_{m}(\theta)\binom{x}{\xi}+U_{m}(\theta), \quad E:=\left[\begin{array}{cc}
0 & -\mathbf{1}  \tag{3.36}\\
\mathbf{1} & 0
\end{array}\right]
$$

where $B_{m}(\theta)$ is a skew-symmetric matrix of dimension $2 d \times 2 d$ and $U_{m}(\theta) \in \mathbb{R}^{2 d}$, and they are both of size $\epsilon_{m}$. Then $\Psi_{m}$ written in the real variables has the form

$$
\begin{equation*}
\Psi_{m}(\theta, x, \xi)=e^{B_{m}(\theta)}(x, \xi)+T_{m}(\theta), \quad \text { where } T_{m}(\theta):=\int_{0}^{1} e^{(1-s) J B_{m}(\theta)} U_{m}(\theta) d s \tag{3.37}
\end{equation*}
$$

Lemma 3.5. There exists a sequence of Hamiltonian matrices $A_{l}(\theta)$ and vectors $V_{l}(\theta) \in \mathbb{R}^{2 d}$ such that

$$
\begin{equation*}
\Psi_{1} \circ \ldots \circ \Psi_{l}(x, \xi)=e^{A_{l}(\theta)}(x, \xi)+V_{l}(\theta) \quad \forall(x, \xi) \in \mathbb{R}^{2 d} \tag{3.38}
\end{equation*}
$$

Furthermore, there exist an Hamiltonian matrix $A_{\omega}(\theta)$ and a vector $V_{\omega}(\theta) \in \mathbb{R}^{2 d}$ such that

$$
\begin{array}{r}
\lim _{l \rightarrow+\infty} e^{A_{l}(\theta)}=e^{A_{\infty}(\theta)}, \\
\lim _{l \rightarrow+\infty} V_{l}(\theta)=V_{\infty}(\theta)  \tag{3.39}\\
|\operatorname{sm} \theta| \leq \sigma / 2
\end{array}\left\|A_{\omega}(\theta)\right\| \leq C \epsilon, \quad \sup _{|\operatorname{Im} \theta| \leq \sigma / 2}\left|V_{\omega}(\theta)\right| \leq C \epsilon,
$$

and for each $\theta \in \mathbb{T}^{n}$,

$$
\Psi(\theta, x, \xi)=e^{A_{\omega}(\theta)}(x, \xi)+V_{\omega}(\theta) \quad \forall(x, \xi) \in \mathbb{R}^{2 d}
$$

Proof. Recall that $\phi_{j}=e^{B_{j}}+T_{j}$ where $T_{j}$ is a translation by the vector $T_{j}$ with the estimates $\left\|B_{j}\right\| \leq C \epsilon_{j},\left\|T_{j}\right\| \leq C \epsilon_{j}$. So we have $e^{B_{j}}=\mathbb{I}+S_{j}$ with $\left\|S_{j}\right\| \leq C \epsilon_{j}$. Then the infinite product $\prod_{1 \leq j<+\infty} e^{B_{j}}$ is convergent. Moreover we have $\prod_{1 \leq j \leq l} e^{B_{j}}=\mathbb{I}+M_{l}$ with $\left\|M_{l}\right\| \leq C \epsilon$ so we have $\prod_{1 \leq j<+\infty} e^{B_{j}}=\mathbb{I}+M$ with $\|M\| \leq C \epsilon$. This is proved by using

$$
\prod_{1 \leq j \leq l}\left(\mathbb{I}+S_{j}\right)=\mathbb{I}+S_{l}+S_{l-1} S_{l}+\cdots S_{1} S_{2} \cdots S_{l}
$$

and estimates on $\left\|S_{j}\right\|$.
So, $M_{l}$ has a small norm and therefore $A_{l}:=\log \left(\mathbb{I}+M_{l}\right)$ is well defined. Furthermore, by construction $\mathbb{I}+M_{l} \in \operatorname{Sp}(2 d)$ and therefore its logarithm is a Hamiltonian matrix, namely $A_{l} \in \operatorname{sp}(2 d)$ for $1 \leq l \leq+\infty$.
Now we have to include the translations. By induction on $l$ we have

$$
\phi_{1} \circ \ldots \circ \phi_{l}(x, \xi)=\mathrm{e}^{A_{l}}(x, \xi)+V_{l},
$$

with $V_{l+1}=\mathrm{e}^{A_{l}} T_{l+1}+V_{l}$ and $V_{1}=T_{1}$. Using the previous estimates we have

$$
\left\|V_{l+1}-V_{l}\right\| \leq C\left\|T_{l+1}\right\| \leq C \epsilon_{l}
$$

Then we get that $\lim _{l \rightarrow+\infty} V_{l}=V_{\infty}$ exists.

## A An example of growth of Sobolev norms (following Graffi and Yajima)

In this appendix we are going to study the Hamiltonian

$$
\begin{equation*}
H:=-\frac{1}{2} \partial_{x x}+\frac{x^{2}}{2}+a x \sin \omega t \tag{A.1}
\end{equation*}
$$

and prove that it is reducible to the Harmonic oscillator if $\omega \neq \pm 1$, while the system exhibits growth of Sobolev norms in the case $\omega= \pm 1$. Actually the result holds in a quite more general situation, but we think that the present example can give a full understanding of the situation with as little techniques as possible. We also remark that in this case it is not necessary to assume that the time dependent part is small.

Finally we recall that A.1 with $\omega= \pm 1$ was studied by Graffi and Yajima as an example of Hamiltonian whose Floquet spectrum is absolutely continuous (despite the fact that the unperturbed Hamiltonian has discrete spectrum). Exploiting the results of EV83, BJL ${ }^{+}$91 one can conclude from [GY00 that the expectation value of the energy is not bounded in this model. The novelty of the present result rests in the much more precise statement ensuring growth of Sobolev norms.

As we already pointed out, in order to get reducibility of the Hamiltonian A.1), it is enough to study the corresponding classical Hamiltonian, in particular proving its reducibility; this is what we will do. It also turns out that all the procedure is clearer working as much as possible at the level of the equations.

So, consider the classical Hamiltonian system

$$
\begin{equation*}
h:=\frac{x^{2}+\xi^{2}}{2}+a x \sin (\omega t), \tag{A.2}
\end{equation*}
$$

whose equations of motion are

$$
\left\{\begin{array}{c}
\dot{x}=\xi  \tag{A.3}\\
\dot{\xi}=-x-a \sin (\omega t)
\end{array} \Longleftrightarrow \quad \ddot{x}+x+a \sin (\omega t)=0 .\right.
$$

Proposition A.1. Assume that $\omega \neq \pm 1$. Then there exists a time periodic canonical transformation conjugating A.2 to

$$
\begin{equation*}
h^{\prime}:=\frac{x^{2}+\xi^{2}}{2} . \tag{A.4}
\end{equation*}
$$

If $\omega= \pm 1$ then the system is canonically conjugated to

$$
\begin{equation*}
h^{\prime}:= \pm \frac{a}{2} \xi \tag{A.5}
\end{equation*}
$$

In both cases the transformation has the form 2.10.
Corollary A.2. In the case $\omega= \pm 1$, for any $s>0$ and $\psi_{0} \in \mathcal{H}^{s}$, there exists a constant $0<C_{s}=C_{s}\left(\left\|\psi_{0}\right\|_{\mathcal{H}^{s}}\right)$ s.t. the solution of the Schrödinger equation with Hamiltonian (A.1) and initial datum $\psi_{0}$ fulfills

$$
\begin{equation*}
\|\psi(t)\|_{\mathcal{H}^{s}} \geq C_{s}\langle t\rangle^{s}, \quad \forall t \in \mathbb{R} \tag{A.6}
\end{equation*}
$$

Before proving the theorem, recall that by the general result of MR17, Theorem 1.5], any solution of the Schrödinger equation with Hamiltonian A.1 fulfills the a priori bound

$$
\begin{equation*}
\|\psi(t)\|_{\mathcal{H}^{s}} \leq C_{s}^{\prime}\left(\left\|\psi_{0}\right\|_{\mathcal{H}^{s}}+|t|^{s}\left\|\psi_{0}\right\|_{\mathcal{H}^{0}}\right), \forall t \in \mathbb{R} \tag{A.7}
\end{equation*}
$$

which is therefore sharp.
Proof of Proposition A.1. We look for a translation

$$
\begin{equation*}
x=x^{\prime}-f(t), \quad \xi=\xi^{\prime}-g(t) \tag{A.8}
\end{equation*}
$$

with $f$ and $g$ time periodic functions to be determined in such a way to eliminate time from A.3. Writing the equations for $\left(x^{\prime}, \xi^{\prime}\right)$, one gets

$$
\dot{x}^{\prime}=\xi^{\prime}-g+\dot{f}, \quad \dot{\xi}^{\prime}=-x^{\prime}-a \sin (\omega t)+\dot{g}+f,
$$

which reduces to the harmonic oscillator by choosing

$$
\left\{\begin{array}{c}
-a \sin (\omega t)+\dot{g}+f=0  \tag{A.9}\\
-g+\dot{f}=0
\end{array} \Longleftrightarrow \ddot{f}+f=a \sin (\omega t)\right.
$$

which has a solution of period $2 \pi / \omega$ only if $\omega \neq \pm 1$. In such a case the only solution having the correct period is

$$
f=\frac{a}{1-\omega^{2}} \sin (\omega t), \quad g=\frac{a \omega}{1-\omega^{2}} \cos (\omega t)
$$

Then the transformation A.8 is a canonical transformation generated as the time one flow of the auxiliary Hamiltonian

$$
\chi:=-\xi \frac{a}{1-\omega^{2}} \sin (\omega t)+x \frac{a \omega}{1-\omega^{2}} \cos (\omega t)
$$

which thus conjugates the classical Hamiltonian A.2 to the Harmonic oscillator; of course the quantization of $\chi$ conjugates the quantum system to the quantum Harmonic oscillator, as follows by Proposition 2.9 .

We come to the resonant case, and, in order to fix ideas, we take $\omega=1$. In such a case the flow of the Harmonic oscillator is periodic of the same period of the forcing, and thus its flow can be used to reduce the system.

In a slight more abstract way, consider a Hamiltonian system with Hamiltonian

$$
H:=\frac{1}{2}\langle z ; B z\rangle+\langle z ; b(t)\rangle
$$

with $z:=(x, \xi), B$ a symmetric matrix, and $b(t)$ a vector valued time periodic function. Then, using the formula 2.4, it is easy to see that the auxiliary time dependent Hamiltonian

$$
\begin{equation*}
\chi_{1}:=\frac{t}{2}\langle z ; B z\rangle \tag{A.10}
\end{equation*}
$$

generates a time periodic transformation which conjugates the system to

$$
h^{\prime}:=\left\langle z ; e^{-J B t} b(t)\right\rangle
$$

( $J$ being the standard symplectic matrix). An explicit computation shows that in our case

$$
\begin{equation*}
h^{\prime}=\frac{a}{2} x \sin (2 t)-\frac{a}{2} \xi \cos (2 t)+\frac{a}{2} \xi \tag{A.11}
\end{equation*}
$$

Then in order to eliminate the two time periodic terms in A.11 it is sufficient to use the canonical transformation generated by the Hamiltonian

$$
\begin{equation*}
\chi_{2}:=-\xi \frac{a}{4} \sin (2 t)-x \frac{a}{4} \cos (2 t) \tag{A.12}
\end{equation*}
$$

which reduce to A.5.
Proof of Corollary A.2. To fix ideas we take $\omega=1$. Let $\chi_{1}^{w} \equiv \frac{t}{2}\left(-\partial_{x x}+x^{2}\right)$ and $\chi_{2}^{w}$ be the Weyl quantization of the Hamiltonians A.10 respectively A.12. By the proof of Proposition A.1. the changes of coordinates

$$
\begin{equation*}
\psi=e^{-\mathrm{i} t H_{0}} \psi_{1}, \quad \psi_{1}=e^{-\mathrm{i} \chi_{2}^{w}(t, x, D)} \varphi, \quad H_{0}:=\frac{1}{2}\left(-\partial_{x x}+x^{2}\right) \tag{A.13}
\end{equation*}
$$

conjugate the Schrödinger equation with Hamiltonian A.1 to the Schrödinger equation with Hamiltonian A.2, namely the transport equation

$$
\partial_{t} \varphi=-\frac{a}{2} \partial_{x} \varphi
$$

The solution of this transport equation is given clearly by

$$
\varphi(t, x)=\varphi_{0}\left(x-\frac{a}{2} t\right)
$$

where $\varphi_{0}$ is the initial datum. Now a simple computation shows that

$$
\liminf _{|t| \rightarrow+\infty}|t|^{-s}\|\varphi(t)\|_{\mathcal{H}^{s}} \geq\left(\frac{|a|}{2}\right)^{s}\left\|\varphi_{0}\right\|_{\mathcal{H}^{0}}
$$

In particular there exists a constant $0<C_{s}=C_{s}\left(\left\|\varphi_{0}\right\|_{\mathcal{H}^{s}}\right)$ s.t.

$$
\begin{equation*}
\|\varphi(t)\|_{\mathcal{H}^{s}} \geq C_{s}\langle t\rangle^{s} \tag{A.14}
\end{equation*}
$$

Since the transformation A.13 maps $\mathcal{H}^{s}$ to $\mathcal{H}^{s}$ uniformly in time (see also Lemma 2.8) estimate A.14 holds also for the original variables.

We remark that by a similar procedure one can also prove the following slightly more general result.

Theorem A.3. Consider the classical Hamiltonian system

$$
\begin{equation*}
h=\sum_{j=1}^{d} \nu_{j} \frac{x_{j}^{2}+\xi_{j}^{2}}{2}+\sum_{j=1}^{d}\left(g_{j}(\omega t) x_{j}+f_{j}(\omega t) \xi_{j}\right) \tag{A.15}
\end{equation*}
$$

with $f_{j}, g_{j} \in C^{r}\left(\mathbb{T}^{n}\right)$.
(1) If there exist $\gamma>0$ and $\tau>n+1$ s.t.

$$
\begin{equation*}
\left|\omega \cdot k \pm \nu_{j}\right| \geq \frac{\gamma}{1+|k|^{\tau}}, \quad \forall k \in \mathbb{Z}^{n}, \quad j=1, \ldots, d \tag{A.16}
\end{equation*}
$$

and $r>\tau+1+n / 2$, then there exists a time quasiperiodic canonical transformation of the form 2.10 conjugating the system $t \$^{5}$

$$
h=\sum_{j=1}^{d} \nu_{j} \frac{x_{j}^{2}+\xi_{j}^{2}}{2}
$$

(2) If there exist $0 \neq \bar{k} \in \mathbb{Z}^{n}$ and $\bar{j}$, s.t.

$$
\begin{equation*}
\omega \cdot \bar{k}-\nu_{\bar{j}}=0 \tag{A.17}
\end{equation*}
$$

and there exist $\gamma>0$ and $\tau$ s.t.

$$
\begin{equation*}
\left|\omega \cdot k \pm \nu_{j}\right| \geq \frac{\gamma}{1+|k|^{\tau}}, \quad \forall(k, j) \neq(\bar{k}, \bar{j}) \tag{A.18}
\end{equation*}
$$

and $r>\tau+1+\frac{n}{2}$, then there exists a time quasiperiodic canonical transformation of the form 2.10 conjugating the system to

$$
h=\sum_{j \neq \bar{j}} \nu_{j} \frac{x_{j}^{2}+\xi_{j}^{2}}{2}+c_{1} x_{\bar{j}}+c_{2} \xi_{\bar{j}},
$$

with $c_{1}, c_{2} \in \mathbb{R}$.
Remark A.4. The constants $c_{1}, c_{2}$ can be easily computed. If at least one of them is different form zero then the solution of the corresponding quantum system exhibits growth of Sobolev norms as in the special model A.1. Of course the result extends in a trivial way to the case in which more resonances are present.

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[^1]:    ${ }^{1}$ We would like to point out also the results PP12 PP15 which at present refer to the resonant nonlinear Schrödinger equation: it would be interesting to study if they have some consequences for reducibility theory.

[^2]:    ${ }^{2}$ recall that a real $2 d \times 2 d$ matrix $A$ belongs to $\mathrm{sp}(2 d)$ iff $J A$ is symmetric

[^3]:    ${ }^{3}$ Formally we could expect $q_{+}$to be of size $O\left(\epsilon^{2}\right)$ but the small divisors and the reduction of the analyticity domain will lead to an estimate of the type $\mathcal{O}\left(\epsilon^{\frac{3}{2}}\right)$.

[^4]:    ${ }^{4}$ i.e. $\left.t \mapsto \lambda_{k}\left(\omega+t z_{k}\right)\right)$ is a holomorphic curve on a neighborhood of 0 , we denote $\partial_{\omega} \lambda(\omega) \cdot z_{k}$ its derivative at $t=0$.

[^5]:    ${ }^{5}$ Actually the transformation is just a translation, so in this case one has $A \equiv 0$.

