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ANALYSIS \& PDE \\ Jean-Marc Delort \\ PERIODIC SOLUTIONS OF NONLINEAR SCHRÖDINGER EQUATIONS: \\ A PARADIFFERENTIAL APPROACH
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# PERIODIC SOLUTIONS OF NONLINEAR SCHRÖDINGER EQUATIONS: A PARADIFFERENTIAL APPROACH 

Jean-Marc Delort


#### Abstract

This paper is devoted to the construction of periodic solutions of nonlinear Schrödinger equations on the torus, for a large set of frequencies. Usual proofs of such results rely on the use of Nash-Moser methods. Our approach avoids this, exploiting the possibility of reducing, through paradifferential conjugation, the equation under study to an equivalent form for which periodic solutions may be constructed by a classical iteration scheme.


## Introduction

This paper is devoted to the existence of families of periodic solutions of Hamiltonian nonlinear Schrödinger equations on the torus $\mathbb{T}^{d}$. Our goal is to show that such results may be proved without using Nash-Moser methods, replacing them by a technically simpler conjugation idea.

We consider equations of type

$$
\left(-i \partial_{t}-\Delta+\mu\right) u=\epsilon \frac{\partial F}{\partial \bar{u}}(\omega t, x, u, \bar{u}, \epsilon)+\epsilon f(\omega t, x)
$$

where $t \in \mathbb{R}, x \in \mathbb{T}^{d}, F$ is a smooth function, vanishing at order 3 at $(u, \bar{u})=0, f$ is a smooth function on $\mathbb{R} \times \mathbb{T}^{d}, 2 \pi$-periodic in time, $\omega$ a frequency parameter, $\mu$ a real parameter and $\epsilon>0$ a small number. One wants to show that for $\epsilon$ small and $\omega$ in a Cantor set whose complement has small measure, the equation has time periodic solutions.

Let us recall known results for that type of problems. The first periodic solutions for nonlinear wave or Schrödinger equations were constructed in [Kuksin 1993; Wayne 1990], which deal with one space dimension, with $x$ staying in a compact interval, and imposing on the extremities of this interval convenient boundary conditions. Later on, Craig and Wayne [1993; 1994] treated the same problem for time-periodic solutions defined on $\mathbb{R} \times \mathbb{S}^{1}$. Periodic solutions of nonlinear wave equations in higher space dimensions (on $\mathbb{R} \times \mathbb{T}^{d}, d \geq 2$ ) were obtained in [Bourgain 1994]. These results concern nonlinearities which are analytic. More recently, some work has been devoted to the same problem when the nonlinearity is a smooth function: Berti and Bolle [2010] have proved in this setting existence of time-periodic solutions for the nonlinear wave equation on $\mathbb{R} \times \mathbb{T}^{d}$. We refer also to the paper of Berti, Bolle and Procesi [Berti et al. 2010], where the case of equations on Zoll manifolds is treated. Very recently, Berti and Procesi

[^0][2011] have studied the same problem, for wave or Schrödinger equations, on a homogeneous space. We refer also to [Craig 2000; Kuksin 2000] for more references.

The proofs of all these results rely on the use of the Nash-Moser theorem, to overcome unavoidable losses of derivatives coming from the small divisors appearing when inverting the linear part of the equation. Our goal here is to show that one may construct periodic solutions of nonlinear Schrödinger equations (for large sets of frequencies), using just a standard iterative scheme instead of the quadratic scheme of the Nash-Moser method. This approach allows one to separate on the one hand the treatment of losses of derivatives coming from small divisors, and on the other hand the question of convergence of the sequence of approximations, while in a Nash-Moser scheme, both problems have to be treated at the same time. The basic idea is inspired by our work in [Delort 2010] concerning linear Schrödinger equations with smooth time dependent potential. It is shown in that paper that a linear equation of type $\left(i \partial_{t}-\Delta+V(t, x)\right) u=0$ may be reduced by conjugation to an equation of type $\left(i \partial_{t}-\Delta+V_{\mathrm{D}}\right) v=R v$, where $R$ is a smoothing operator and $V_{\mathrm{D}}$ a block diagonal operator of order zero. We aim at applying a similar method when the linear potential $V$ is replaced by a nonlinear one, so that, in the reduced equation, the block-diagonal operator $V_{\mathrm{D}}$ depends on $v$ itself, and $R$ sends essentially $H^{s}$ to $H^{2 s-a}$ (where $a$ is a fixed constant, and $H^{s}$ the Sobolev scale). It is pretty clear that such a reduced equation will be solvable by a standard iterative scheme, even if the inversion of $i \partial_{t}-\Delta+V_{\mathrm{D}}$ loses derivatives because of small divisors, since such losses are recovered by the smoothing properties of $R$ on the right side.

Before describing the different sections of the paper, let us give some more references and add some comments. There are actually a few results concerning existence of periodic solutions which do not appeal to Nash-Moser theorem. Bambusi and Paleari [2001; 2002] constructed such solutions without making use of Nash-Moser or KAM methods, but only for a family of frequency parameters of measure zero (instead of a set of parameters whose complement has small measure). Related results, concerning the case of rational frequencies, may be found in [Berti 2007, Chapter 5]. Recently, Gentile and Procesi [2009] found, for analytic nonlinearities, an alternative approach to Nash-Moser using expansions in terms of Lindsted series.

Let us also mention that we restrict in this paper to one of the many variants that may be considered when constructing periodic solutions. Most of the known results we cited so far concern the case of periodic solutions of the nonlinear equation, whose frequency is close to the frequency of a periodic solution of the linear equation obtained for $\epsilon=0$. The problem may be written, using a LiapunovSchmidt decomposition, as a coupling between a non-resonant equation (the ( $P$ ) equation) and a resonant one (the $(Q)$ equation). In most works, the resonant equation is a finite-dimensional equation, while $(P)$ is infinite-dimensional. One uses Nash-Moser to solve $(P)$, getting a solution depending on finitely many parameters. Plugging this solution in $(Q)$, one gets for these finitely many parameters an equation in closed form, that may be solved using implicit functions-like theorems. Actually, Berti and Bolle [2006] have shown that such a strategy may be also adapted to the case when $(Q)$ is completely resonant, i.e., infinite-dimensional.

Since our objective here is to show that one may avoid the use of Nash-Moser theorems, we limited ourselves to the forced oscillations equation written at the beginning of the introduction, which
corresponds to a $(P)$ equation for which there is no associated ( $Q$ ) equation. Berti and Bolle [2010] have studied similar forced oscillations for the wave equation. It is very likely that our method could be adapted to recover as well known results for resonant periodic Schrödinger equations, even if one would have to write a detailed proof. In the same way, since the results in [Delort 2010] concerning the Schrödinger equation hold not only on $\mathbb{T}^{d}$, but also on Zoll manifolds or on some surfaces of revolution, we conjecture that the analogue of the main theorem of this paper extends to this setting, or even to the case of a product of several Zoll manifolds.

Organization of the paper. Section 1 states the main theorem and introduces notation.
Section 2 is devoted to the paralinearization of the equation. After defining convenient classes of paradifferential operators, we perform a first reduction, localizing the unknown of the problem close to the characteristic variety of the linear Schrödinger operator. This is done using the standard implicit function theorem. Next, we paralinearize the equation, reducing it to

$$
\left(-i \omega \partial_{t}-\Delta+V\right) v=R(v) v+\epsilon f
$$

where $V$ is a paradifferential operator of order zero, depending on $v$, and $R(v)$ is a smoothing operator (Actually, we shall have to consider a system in $(v, \bar{v})$ instead of a scalar equation). A consequence of the fact that our starting equation is Hamiltonian will be that $V$ is self-adjoint.

Section 3 is the heart of the paper. We construct a paradifferential conjugation of the preceding equation to transform it into

$$
\left(-i \omega \partial_{t}-\Delta+V_{\mathrm{D}}(w)\right) w=R(w) w+\epsilon f
$$

where $R(w)$ is still a smoothing operator, and $V_{\mathrm{D}}$ is block diagonal relatively to an orthogonal decomposition of $L^{2}\left(\mathbb{T}^{d}\right)$ in a sum of finite-dimensional subspaces introduced in [Bourgain 1999].

Section 4 is devoted to the construction of the solution to the block diagonal equation by a standard iteration scheme. We first show that on each block $-i \omega \partial_{t}-\Delta+V_{\mathrm{D}}(w)$ is invertible for $\omega$ outside a convenient small subset. This is done by the usual argument, exploiting that the $\omega$-derivative of the eigenvalues of $-i \omega \partial_{t}-\Delta$ is large. In order that the set of excluded parameters remain small, we have to allow small divisors when inverting $-i \omega \partial_{t}-\Delta+V_{\mathrm{D}}(w)$. As the right-hand side of the equation involves a smoothing operator $R(w)$, we may compensate the losses of derivatives coming from such small divisors, and construct a sequence of approximations of the solution.

Let us conclude this introduction with a few words concerning the limitations of our method. First, it does not seem that it could be adapted to find periodic solutions of nonlinear wave equations, as the construction of Section 3 relies on a specific separation property for the eigenvalues of $-\Delta$ on $\mathbb{T}^{d}$. On the other hand, it might be applied to equations where one has a nice separation of eigenvalues, like KdV or one-dimensional water wave equations with surface tension. Second, we do not know if our method could be modified to construct quasi-periodic solutions. Recall that such solutions have been obtained for the equation set on an interval [Kuksin 1993; Kuksin 2000; Kuksin and Pöschel 1996]. The case of solutions on $\mathbb{S}^{1}$ has been treated in [Bourgain 1994]. In higher dimensions, Bourgain [1998] constructed such periodic solutions on $\mathbb{T}^{2}$. The case of general $\mathbb{T}^{d}$ has been treated in [Bourgain 2005;

Eliasson and Kuksin 2010]. One of the difficulties of the quasi-periodic case versus the periodic one lies in the fact that, even close to the characteristic variety, time frequencies might be much larger than space frequencies. In our proof below, the fact that these frequencies are of the same magnitude plays an important role. We do not know whether the multiscale methods of Bourgain, Eliasson, and Kuksin could be combined to the arguments we use in the periodic case to construct quasi-periodic solutions without making appeal to a Newton scheme.

## 1. Periodic solutions of semi-linear Schrödinger equations

1.1. Statement of the main theorem. Let $\mathbb{T}^{d}(d \geq 1)$ be the standard torus, $\mathbb{S}^{1}$ the unit circle. Consider a $C^{\infty}$ function

$$
\begin{align*}
& F:(t, x, u, \bar{u}, \epsilon) \longrightarrow F(t, x, u, \bar{u}, \epsilon) \\
& \mathbb{R} \times \mathbb{T}^{d} \times \mathbb{C}^{2} \times[0,1] \rightarrow \mathbb{R} \tag{1.1.1}
\end{align*}
$$

which is $2 \pi$-periodic in $t$, and satisfies $\partial_{u, \bar{u}}^{\alpha} F(t, x, 0,0, \epsilon) \equiv 0$ for $|\alpha| \leq 2$. We study the equation

$$
\begin{equation*}
\left(D_{t}-\Delta+\mu\right) u=\epsilon \frac{\partial F}{\partial \bar{u}}(\omega t, x, u, \bar{u}, \epsilon)+\epsilon f(\omega t, x) \tag{1.1.2}
\end{equation*}
$$

where $\Delta$ is the Laplace operator on $\mathbb{T}^{d}, D_{t}=\frac{1}{i} \frac{\partial}{\partial t}, \epsilon \in[0,1], \mu \in \mathbb{R}, \omega \in \mathbb{R}_{+}^{*}, f$ is a smooth function on $\mathbb{R} \times \mathbb{T}^{d}, 2 \pi$-periodic in $t$, with values in $\mathbb{C}$, and where we look for $\frac{2 \pi}{\omega}$-periodic solutions of the equation when $\epsilon$ is small. Changing $t$ to $t / \omega$, we have to find solutions on $\mathbb{S}^{1} \times \mathbb{T}^{d}$ to the equivalent equation

$$
\begin{equation*}
\left(\omega D_{t}-\Delta+\mu\right) u=\epsilon \frac{\partial F}{\partial \bar{u}}(t, x, u, \bar{u}, \epsilon)+\epsilon f(t, x) \tag{1.1.3}
\end{equation*}
$$

for small enough $\epsilon$ and for $\omega$ outside a subset of small measure. To fix ideas, we shall take $\omega$ inside a fixed compact subinterval of $] 0,+\infty[$, say $\omega \in[1,2]$.

Let us define the Sobolev space in which we shall look for solutions. If $u \in \mathscr{D}^{\prime}\left(\mathbb{S}^{1} \times \mathbb{T}^{d}\right)$, we set for $(j, n) \in \mathbb{Z} \times \mathbb{Z}^{d}$

$$
\hat{u}(j, n)=\frac{1}{(2 \pi)^{\frac{d+1}{2}}} \int_{\mathbb{S}^{1} \times \mathbb{T}^{d}} e^{-i t j-i n \cdot x} u(t, x) d t d x
$$

For $s \in \mathbb{R}$, define $\widetilde{\mathscr{H}}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}\right)$ to be the space of those $u \in \mathscr{D}^{\prime}\left(\mathbb{S}^{1} \times \mathbb{T}^{d}\right)$ such that

$$
\begin{equation*}
\|u\|_{\tilde{\mathscr{H}}}^{2} \stackrel{\text { def }}{=} \sum_{j \in \mathbb{Z}} \sum_{n \in \mathbb{Z}^{d}}\left(1+|j|+|n|^{2}\right)^{s}|\hat{u}(j, n)|^{2}<+\infty \tag{1.1.4}
\end{equation*}
$$

We shall use the similar notation $\tilde{\mathscr{H}}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}^{2}\right), \tilde{\mathscr{H}}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right)$ for $\mathbb{C}^{2}$ or $\mathbb{R}^{2}$-valued functions.
Let us state our main theorem.
Theorem 1.1.1. Let $\mu \in \mathbb{R}-\mathbb{Z}_{-}$. There are $s_{0}>0, \zeta>0$ and for any $s \geq s_{0}$, any $q_{0}>0$, there are constants $\left.\left.\delta_{0} \in\right] 0,1\right], B>0$ and for any $f \in \tilde{\mathscr{H}}^{s+\zeta}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}\right)$ with $\|f\|_{\tilde{H}^{s+\zeta}} \leq q_{0}$, there is a subset $\bigcirc \subset[1,2] \times] 0,1]$ such that:

- For any $\left.\delta \in] 0, \delta_{0}\right]$ and $\epsilon \in\left[0, \delta^{2}\right]$

$$
\begin{equation*}
\operatorname{meas}\{\omega \in[1,2]:(\omega, \epsilon) \in \mathcal{O}\} \leq B \delta \tag{1.1.5}
\end{equation*}
$$

- For any $\left.\delta \in] 0, \delta_{0}\right]$, any $\epsilon \in\left[0, \delta^{2}\right]$, and any $\omega \in[1,2]$ such that $(\omega, \epsilon) \notin \mathcal{O}$, (1.1.3) has a solution $u \in \tilde{\mathscr{H}}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}\right)$ satisfying $\|u\|_{\mathscr{H}^{s}} \leq B \in \delta^{-1}$.
Remark. As mentioned in the introduction, this theorem is a version, for Schrödinger equations, of [Berti and Bolle 2010, Theorem 1.1], which concerns wave equations. Our point will be to give a proof that does not appeal to Nash-Moser methods.
1.2. Spaces offunctions and notations. For $n \in \mathbb{Z}^{d}, u \in \mathscr{D}^{\prime}\left(\mathbb{T}^{d}\right)$, we denote by $\Pi_{n}$ the spectral projector

$$
\begin{equation*}
\Pi_{n} u=\hat{u}(n) \frac{e^{i n \cdot x}}{(2 \pi)^{d / 2}}=\int_{\mathbb{T}^{d}} e^{-i n \cdot x} u(x) \frac{d x}{(2 \pi)^{d / 2}} \frac{e^{i n \cdot x}}{(2 \pi)^{d / 2}} \tag{1.2.1}
\end{equation*}
$$

When $u(t, x)$ is in $\mathscr{D}^{\prime}\left(\mathbb{S}^{1} \times \mathbb{T}^{d}\right)$, we use the same notation, considering $t$ as a parameter. We shall make use of the following "separation property" result attributed by Bourgain to Granville and Spencer (see [Bourgain 1999, Lemma 8.1]; for the proof see also [Bourgain 2005, Lemma 19.10]).
Lemma 1.2.1. For any $\beta \in] 0, \frac{1}{10}[$, there are $\rho \in] 0, \beta\left[, \theta>0\right.$ and a partition $\left(\Omega_{\alpha}\right)_{\alpha \in \mathscr{A}}$ of $\mathbb{Z}^{d}$ such that

$$
\begin{align*}
& \left|n-n^{\prime}\right|+\left||n|^{2}-\left|n^{\prime}\right|^{2}\right|<\theta+|n|^{\beta} \quad \text { for all } \alpha \in \mathscr{A}, n \in \Omega_{\alpha}, n^{\prime} \in \Omega_{\alpha} \\
& \quad\left|n-n^{\prime}\right|+\left||n|^{2}-\left|n^{\prime}\right|^{2}\right|>|n|^{\rho} \quad \text { for all } \alpha, \alpha^{\prime} \in \mathscr{A}\left(\alpha \neq \alpha^{\prime}\right), n \in \Omega_{\alpha}, n^{\prime} \in \Omega_{\alpha^{\prime}} . \tag{1.2.2}
\end{align*}
$$

For each $\alpha \in \mathscr{A}$, we choose some $n(\alpha) \in \Omega_{\alpha}$. There is a constant $\Theta_{0}>0$ such that, if we set $\langle n\rangle=\left(1+|n|^{2}\right)^{1 / 2}$ for $n \in \mathbb{Z}^{d}$, then

$$
\begin{equation*}
\Theta_{0}^{-1}\langle n(\alpha)\rangle \leq\langle n\rangle \leq \Theta_{0}\langle n(\alpha)\rangle \tag{1.2.3}
\end{equation*}
$$

for any $\alpha \in \mathscr{A}$, any $n \in \Omega_{\alpha}$. It also follows from (1.2.2) that, for some uniform constant $\Theta_{1}>0$,

$$
\begin{equation*}
\# \Omega_{\alpha} \leq \Theta_{1}\langle n(\alpha)\rangle^{\beta d} \tag{1.2.4}
\end{equation*}
$$

For any $\alpha \in \mathscr{A}$, we set

$$
\begin{equation*}
\tilde{\Pi}_{\alpha}=\sum_{n \in \Omega_{\alpha}} \Pi_{n} \tag{1.2.5}
\end{equation*}
$$

We define a closed subspace $\mathscr{H}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}\right)$ of $\tilde{\mathscr{H}}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}\right)$ by

$$
\begin{array}{r}
\mathscr{H}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}\right)=\bigcap_{\alpha \in \mathscr{A}}\left\{u \in \tilde{\mathscr{H}}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}\right): \hat{u}(j, n)=0 \text { for all } n \in \Omega_{\alpha} \text { and all } j\right. \text { such that }  \tag{1.2.6}\\
\left.|j|>K_{0}\langle n(\alpha)\rangle^{2} \text { or }|j|<K_{0}^{-1}\langle n(\alpha)\rangle^{2}\right\},
\end{array}
$$

where $K_{0}=K_{0}(\mu)$ will be chosen later on.
In other words, non vanishing modes $(j, n)$ of an element $u$ of $\mathscr{H}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}\right)$ have to satisfy $K_{0}^{-1}\langle n(\alpha)\rangle^{2} \leq|j| \leq K_{0}\langle n(\alpha)\rangle^{2}$ if $n \in \Omega_{\alpha}$. This shows that the restriction to $\mathscr{H}^{s}$ of the $\tilde{\mathscr{H}}^{s}$-norm
given by (1.1.4) is equivalent to the square root of

$$
\begin{equation*}
\sum_{j \in \mathbb{Z}} \sum_{n \in \mathbb{Z}^{d}}\langle n\rangle^{2 s}|\hat{u}(j, n)|^{2} \tag{1.2.7}
\end{equation*}
$$

and to the square root of

$$
\begin{equation*}
\sum_{\alpha \in \mathscr{A}}\langle n(\alpha)\rangle^{2 s}\left\|\tilde{\Pi}_{\alpha} u\right\|_{L^{2}\left(\mathbb{S}^{1} \times \mathbb{T}^{d}, \mathbb{C}\right)}^{2} \tag{1.2.8}
\end{equation*}
$$

We use similar notation for spaces $\mathscr{H}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}^{2}\right), \mathscr{H}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right)$, and so on.

## 2. Paralinearization of the equation

The goal of this section is to rewrite (1.1.3) as a paradifferential equation in the sense of [Bony 1981], on spaces of form (1.2.6). We first define the classes of operators we shall use.
2.1. Spaces of operators. We fix from now on some real number $\sigma_{0}>\frac{d}{2}+1$. If $s \in \mathbb{R}, q>0$, we denote by $B_{q}\left(\mathscr{H}^{s}\right)$ the open ball with center 0 , radius $q$ in $\mathscr{H}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}\right)$, $\mathscr{H}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}^{2}\right), \ldots$
Definition 2.1.1. Let $m \in \mathbb{R}, q>0, N \in \mathbb{N}, \sigma \in \mathbb{R}, \sigma \geq \sigma_{0}+2 N+d+1$. One denotes by $\Psi^{m}(N, \sigma, q)$ the space of maps $U \rightarrow a(U)$ defined on the open ball of center 0 , radius $q$ in $\mathscr{H}^{\sigma}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}^{2}\right)$, with values in the space of linear maps from $C^{\infty}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}\right)$ to $\mathscr{D}^{\prime}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}\right)$, such that, for any $n, n^{\prime} \in \mathbb{Z}^{d}$, the map $U \rightarrow \Pi_{n} a(U) \Pi_{n^{\prime}}$ is smooth with values in $\mathscr{L}\left(\mathscr{H}^{0}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}\right)\right)$ and satisfies for any $M \in \mathbb{N}$ with $d+1 \leq M \leq \sigma-\sigma_{0}-2 N$, any $U \in B_{q}\left(\mathscr{H}^{\sigma}\right)$, any $j \in \mathbb{N}$, any $W_{1}, \ldots, W_{j} \in \mathscr{H}^{\sigma}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}^{2}\right)$, any $n, n^{\prime} \in \mathbb{Z}^{d}$,

$$
\begin{align*}
\| \Pi_{n}\left(\partial _ { U } ^ { j } a ( U ) \cdot \left(W_{1},\right.\right. & \left.\left.\ldots, W_{j}\right)\right) \Pi_{n^{\prime}} \|_{\mathscr{L}\left(\mathscr{H}^{0}\right)} \\
& \leq C\left(1+|n|+\left|n^{\prime}\right|\right)^{m}\left\langle n-n^{\prime}\right\rangle^{-M_{1}} \mathbb{1}_{\left|n-n^{\prime}\right| \leq \frac{1}{10}}\left(|n|+\left|n^{\prime}\right|\right) \tag{2.1.1}
\end{align*} \prod_{\ell=1}^{j}\left\|W_{\ell}\right\|_{\mathscr{H} \sigma_{0}+2 N+M} .
$$

Remarks. - In (2.1.1), the decay $\left\langle n-n^{\prime}\right\rangle^{-M}$ reflects the available $x$-smoothness of the symbol of a pseudo-differential or paradifferential operator. This smoothness is controlled by the upper bound $\sigma-\sigma_{0}-2 N$ that we assume for $M$. The cut-off $\left|n-n^{\prime}\right| \leq \frac{1}{10}\left(|n|+\left|n^{\prime}\right|\right)$ means that we are considering paradifferential operators. The integer $N$ measures some loss of smoothness, relatively to the index $\sigma$, that will appear in some expansions of operators.

- Definition 2.1.1 implies that if $a \in \Psi^{m}(N, \sigma, q)$, then $\partial_{t}[a(U)]$ belongs to $\Psi^{m}(N+1, \sigma, q)$. Actually, $\partial_{t} a(U)=\partial_{U} a(U) \cdot \partial_{t} U$, so (2.1.1) allows us to estimate

$$
\left\|\Pi_{n}\left(\partial_{U}^{j}\left(\partial_{t}[a(U)]\right) \cdot\left(W_{1}, \ldots, W_{j}\right)\right) \Pi_{n^{\prime}}\right\|_{\mathscr{L}\left(\mathscr{H}^{0}\right)}
$$

from $\left\|\partial_{t} U\right\|_{\mathscr{H} \sigma_{0}+2 N+M} \prod_{\ell=1}^{j}\left\|W_{\ell}\right\|_{\mathscr{H} \sigma_{0}+2 N+M}$, and by the definition (1.2.6) of $\mathscr{H}^{s}$,

$$
\left\|\partial_{t} U\right\|_{\mathscr{H} \sigma_{0}+2 N+M} \leq K_{0}\|U\|_{\mathscr{H} \sigma_{0}+2(N+1)+M} \leq K_{0}\|U\|_{\mathscr{H} \sigma}
$$

if we assume $M \leq \sigma-2(N+1)-\sigma_{0}$.
The definition implies boundedness properties for the operators.

Lemma 2.1.2. Let $\sigma, m, N, q$ be as in the definition. Assume that $\sigma \geq \sigma_{0}+2 N+d+1$. Then for any $U \in B_{q}\left(\mathscr{H}^{\sigma}\right)$, for any $s \in \mathbb{R}, a(U)$ is a bounded operator from $\mathscr{H}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}\right)$ to $\mathscr{H}^{s-m}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}\right)$. Moreover, $U \rightarrow a(U)$ is a smooth map from $B_{q}\left(\mathscr{H}^{\sigma}\right)$ to the space $\mathscr{L}\left(\mathscr{H}^{s}, \mathscr{H}^{s-m}\right)$, and for any $j \in \mathbb{N}$, there is $C>0$, such that, for any $U \in B_{q}\left(\mathscr{H}^{\sigma}\right)$ and any $W_{1}, \ldots, W_{j} \in \mathscr{H}^{\sigma}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}\right)$,

$$
\begin{equation*}
\left\|\partial_{U}^{j} a(U) \cdot\left(W_{1}, \ldots, W_{j}\right)\right\|_{\mathscr{L}\left(\mathscr{H}^{s}, \mathscr{H}^{s-m}\right)} \leq C \prod_{\ell=1}^{j}\left\|W_{\ell}\right\|_{\mathscr{H} \sigma_{0}+2 N+d+1} \tag{2.1.2}
\end{equation*}
$$

Proof. One has just to apply (2.1.1) with $M=d+1$ and use the fact that, by (1.2.7), $\|v\|_{\mathscr{Y} s}^{2}$ is equivalent to $\sum_{n \in \mathbb{Z}^{d}}\langle n\rangle^{2 s}\left\|\Pi_{n} v\right\|_{L^{2}}^{2}$.

We define a class of smoothing operators as well.
Definition 2.1.3. Let $\sigma \in \mathbb{R}, N \in \mathbb{N}$, and $v \in \mathbb{N}$, with $\sigma \geq \sigma_{0}+2 N+d+1, q>0, r \in \mathbb{R}_{+}$. One denotes by $\mathscr{R}_{\nu}^{r}(N, \sigma, q)$ the space of smooth maps $U \rightarrow R(U)$ defined on $B_{q}\left(\mathscr{H}^{\sigma}\right)$, with values in $\mathscr{L}\left(\mathscr{H}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}\right), \mathscr{H}^{s+r}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}\right)\right)$ for any $s \geq \sigma_{0}+v$, such that there is for any $j$, any $s \geq \sigma_{0}+v$, a constant $C>0$ with

$$
\begin{equation*}
\left\|\partial_{U}^{j} R(U) \cdot\left(W_{1}, \ldots, W_{j}\right)\right\|_{\mathscr{L}\left(\mathscr{H}^{s}, \mathscr{H}^{s+r}\right)} \leq C \prod_{\ell=1}^{j}\left\|W_{\ell}\right\|_{\mathscr{H} \sigma} \tag{2.1.3}
\end{equation*}
$$

for any $U \in B_{q}\left(\mathscr{H}^{\sigma}\right), W_{1}, \ldots, W_{j} \in \mathscr{H}^{\sigma}$.
Remark. Lemma 2.1.2 shows that if $r \geq 0$ and $\sigma \geq \sigma_{0}+2 N+d+1$, the space $\Psi^{-r}(N, \sigma, q)$ is contained in $\mathscr{R}_{0}^{r}(N, \sigma, q)$.

Proposition 2.1.4. (i) Let $\sigma \geq \sigma_{0}+2 N+d+1, a \in \Psi^{m}(N, \sigma, q)$. Then $a^{*} \in \Psi^{m}(N, \sigma, q)$.
(ii) Let $m_{1}, m_{2} \in \mathbb{R}$. Assume $\sigma \geq \sigma_{0}+2 N+d+1+\left(m_{1}+m_{2}\right)_{+}$. Set

$$
\begin{equation*}
r=\sigma-\sigma_{0}-2 N-(d+1)-\left(m_{1}+m_{2}\right) \geq 0 \tag{2.1.4}
\end{equation*}
$$

If $a \in \Psi^{m_{1}}(N, \sigma, q)$ and $b \in \Psi^{m_{2}}(N, \sigma, q)$, there are $c \in \Psi^{m_{1}+m_{2}}(N, \sigma, q)$ and $R \in \mathscr{R}_{0}^{r}(N, \sigma, q)$ such that

$$
\begin{equation*}
a(U) \circ b(U)=c(U)+R(U) \tag{2.1.5}
\end{equation*}
$$

Proof. Part (i) follows immediately from the definition. For (ii), define

$$
c(U)=\sum_{n} \sum_{n^{\prime}} \Pi_{n}[a(U) \circ b(U)] \Pi_{n^{\prime}} \mathbb{1}_{\left|n-n^{\prime}\right| \leq \frac{1}{10}}\left(|n|+\left|n^{\prime}\right|\right)
$$

To check that (2.1.1) is satisfied by $c$ when $j=0$ we write

$$
\left\|\Pi_{n} c(U) \Pi_{n^{\prime}}\right\|_{\mathscr{L}\left(\mathscr{H}^{0}\right)} \leq \sum_{k}\left\|\Pi_{n} a(U) \Pi_{k}\right\|_{\mathscr{L}\left(\mathscr{H}^{0}\right)}\left\|\Pi_{k} b(U) \Pi_{n^{\prime}}\right\|_{\mathscr{L}\left(\mathscr{H}^{0}\right)}
$$

for $n, n^{\prime}$ with $\left|n-n^{\prime}\right| \leq \frac{1}{10}\left(|n|+\left|n^{\prime}\right|\right)$. Applying (2.1.1) to $a, b$ with $d+1 \leq M \leq \sigma-\sigma_{0}-2 N$, we get the bound

$$
C\left(1+|n|+\left|n^{\prime}\right|\right)^{m_{1}+m_{2}} \sum_{k}\langle n-k\rangle^{-M}\left\langle k-n^{\prime}\right\rangle^{-M} \leq C\left(1+|n|+\left|n^{\prime}\right|\right)^{m_{1}+m_{2}}\left\langle n-n^{\prime}\right\rangle^{-M}
$$

One estimates $\partial_{U}^{j} c(U)$ in the same way.
The remainder $R(U)=a(U) \circ b(U)-c(U)$ will satisfy by definition of $c$ :

$$
\left\|\Pi_{n} R(U) \Pi_{n^{\prime}}\right\|_{\mathscr{L}\left(\mathscr{H}^{0}\right)} \leq \sum_{k}\left\|\Pi_{n} a(U) \Pi_{k}\right\|_{\mathscr{L}\left(\mathscr{H}^{0}\right)}\left\|\Pi_{k} b(U) \Pi_{n^{\prime}}\right\|_{\mathscr{L}\left(\mathscr{H}^{0}\right)} \mathbb{1}_{\left|n-n^{\prime}\right|>\frac{1}{10}\left(|n|+\left|n^{\prime}\right|\right)}
$$

and so will be bounded using (2.1.1) for $a, b$ by

$$
C\left(1+|n|+\left|n^{\prime}\right|\right)^{m_{1}+m_{2}} \sum_{k}\langle n-k\rangle^{-M}\left\langle k-n^{\prime}\right\rangle^{-M_{1}} \mathbb{1}_{|k-n| \leq \frac{1}{10}}(|n|+|k|)_{\left|k-n^{\prime}\right| \leq \frac{1}{10}}\left(\left|n^{\prime}\right|+|k|\right) \mathbb{1}_{\left|n-n^{\prime}\right|>\frac{1}{10}}\left(|n|+\left|n^{\prime}\right|\right)
$$

for any $M$ between $d+1$ and $\sigma-\sigma_{0}-2 N$. Since on the summation, either $|n-k| \geq \frac{1}{2}\left|n-n^{\prime}\right|$ or $\left|n^{\prime}-k\right| \geq \frac{1}{2}\left|n-n^{\prime}\right|$, and $\left|n-n^{\prime}\right| \leq \frac{1}{2}\left(|n|+\left|n^{\prime}\right|\right)$, we get the bound

$$
\left\|\Pi_{n} R(U) \Pi_{n^{\prime}}\right\|_{\mathscr{L}\left(\mathscr{H}^{0}\right)} \leq C\left(1+|n|+\left|n^{\prime}\right|\right)^{m_{1}+m_{2}-M_{\mathbb{1}_{\mid n-n^{\prime}} \left\lvert\, \leq \frac{1}{2}\left(|n|+\left|n^{\prime}\right|\right)\right.}}
$$

for any $M$ between $d+1$ and $\sigma-\sigma_{0}-2 N$. Reasoning as in the proof of Lemma 2.1.2, we obtain that $R(U)$ sends $\mathscr{H}^{s}$ to $\mathscr{H}^{s+r}$ for any $s$ and $r$ given by (2.1.4). The estimates of $\partial_{U}^{j} R(U) \cdot\left(W_{1}, \ldots, W_{j}\right)$ are obtained in the same way.

In the rest of this paper, we shall use several variants of these classes. We shall denote by $\Psi_{\mathbb{R}}^{m}(N, \sigma, q)$ the subspace of $\Psi^{m}(N, \sigma, q)$ made of those operators $a(U)$ sending real-valued functions to real-valued functions, i.e., satisfying $\overline{a(U)}=a(U)$. We define $\mathscr{R}_{v, \mathbb{R}}^{r}(N, \sigma, q)$ from $\mathscr{R}_{\nu}^{r}(N, \sigma, q)$ analogously. We denote by

$$
\Psi^{m}(N, \sigma, q) \otimes \mathcal{M}_{2}(\mathbb{R}) \quad \text { and } \quad \mathscr{R}_{v}^{r}(N, \sigma, q) \otimes \mathcal{M}_{2}(\mathbb{R})
$$

the space of $2 \times 2$ matrices with entries in $\Psi^{m}(N, \sigma, q)$ and in $\mathscr{R}_{v}^{r}(N, \sigma, q)$ respectively. We use similar notation for $\Psi_{\mathbb{R}}^{m}(N, \sigma, q)$ and $\mathscr{R}_{v, \mathbb{R}}^{r}(N, \sigma, q)$.

Finally, we shall consider operators $a(U, \omega, \epsilon), R(U, \omega, \epsilon)$ depending on $(\omega, \epsilon)$ staying in a bounded domain of $\mathbb{R}^{2}$. We say $a(U, \omega, \epsilon)$ is $C^{1}$ in $(\omega, \epsilon)$ if $(\omega, \epsilon) \rightarrow \Pi_{n} a(U, \omega, \epsilon) \Pi_{n^{\prime}}$ is $C^{1}$ in $(\omega, \epsilon)$ with values in $\mathscr{L}\left(\mathscr{H}^{0}\right)$ and if (2.1.1) is satisfied also by $\partial_{\omega} a, \partial_{\epsilon} a$. Likewise, $R(U, \omega, \epsilon)$ is $C^{1}$ if $(\omega, \epsilon) \rightarrow R(U, \omega, \epsilon)$ is $C^{1}$ in $(\omega, \epsilon)$ with values in $\mathscr{L}\left(\mathscr{H}^{s}, \mathscr{H}^{s+r}\right)$ and if (2.1.3) is satisfied by $\partial_{\omega} R, \partial_{\epsilon} R$.
2.2. Equivalent formulation of the equation. The goal of this subsection is to reduce (1.1.3) to an equivalent equation for a new unknown belonging to the space $\mathscr{H}^{s}$ defined by (1.2.6) instead of $\tilde{\mathscr{H}}^{s}$. Recall that we fixed some $\sigma_{0}>\frac{d}{2}+1$.

For $\sigma \in \mathbb{R}$, we consider the space $\mathscr{H}^{\sigma}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right) \subset \tilde{\mathscr{H}}^{\sigma}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right)$ and denote by $\mathscr{F}^{\sigma}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right)$ the orthogonal complement of the first space in the second one.

Definition 2.2.1. Let $\sigma \geq \sigma_{0}$. Denote by $\mathscr{H}_{1}^{\sigma}$, $\mathscr{H}_{2}^{\sigma}$ any of the preceding spaces. Let $X$ be an open subset of $\mathscr{H}_{1}^{\sigma}, k \in \mathbb{Z}$. One denotes by $\Phi^{\infty, k}\left(X, \mathscr{H}_{2}^{\sigma-k}\right)$ the space of $C^{\infty}$ maps $G: X \rightarrow \mathscr{H}_{2}^{\sigma-k}$ such that, for any $s \geq \sigma$ and $u \in X \cap \mathscr{H}_{1}^{s}$ :

- $G(u) \in \mathscr{H}_{2}^{s-k}$.
- The linear map $D G(u) \in \mathscr{L}\left(\mathscr{H}_{1}^{\sigma}, \mathscr{H}_{2}^{\sigma-k}\right)$ extends as an element of $\mathscr{L}\left(\mathscr{H}_{1}^{\sigma^{\prime}}, \mathscr{H}_{2}^{\sigma^{\prime}-k}\right)$ for any $\sigma^{\prime} \in[-s, s]$. Moreover, $v \rightarrow D G(v)$ is smooth from $X \cap \mathscr{H}_{1}^{s}$ to the preceding space.
- The bilinear map $D^{2} G(u) \in \mathscr{L}_{2}\left(\mathscr{H}_{1}^{\sigma} \times \mathscr{H}_{1}^{\sigma} ; \mathscr{H}_{2}^{\sigma-k}\right)$ extends as an element of $\mathscr{L}_{2}\left(\mathscr{H}_{1}^{\sigma_{1}} \times \mathscr{H}_{1}^{\sigma_{2}} ; \mathscr{H}_{2}^{-\sigma_{3}-k}\right)$ for any triple $\left\{\sigma_{1}, \sigma_{2}, \sigma_{3}\right\}=\left\{\sigma^{\prime},-\sigma^{\prime}, \max \left(\sigma_{0}, \sigma^{\prime}\right)\right\}$ with $\sigma^{\prime} \in[0, s]$. Moreover, $v \rightarrow D^{2} G(v)$ is smooth from $X \cap \mathscr{H}_{1}^{s}$ to the preceding space.
Let us give an example of an element of $\Phi^{\infty, 0}\left(\tilde{\mathscr{H}}^{\sigma}, \tilde{\mathscr{H}}^{\sigma}\right)$. Consider $F: \mathbb{S}^{1} \times \mathbb{T}^{d} \times \mathbb{R}^{2} \rightarrow \mathbb{R}^{2}$, a smooth function satisfying $F(t, x, 0) \equiv 0, \partial_{u} F(t, x, 0) \equiv 0$. By Lemma A. 1 of the appendix, for $\sigma>\frac{d}{2}+1$ and $u \in \tilde{\mathscr{H}}^{\sigma}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right)$, we have $F(\cdot, u) \in \tilde{\mathscr{H}}^{\sigma}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right)$, and by Corollary A.2, $u \rightarrow F(\cdot, u)$ is smooth. If we define $G(u)=F(\cdot, u)$, then $D G(u) \cdot h=\partial_{u} F(\cdot, u) h$, which, by Lemma A.3, extends as a linear map from $\tilde{\mathscr{H}}^{\sigma^{\prime}}$ to itself for any $\sigma^{\prime} \in[-s, s]$, when $u \in \widetilde{\mathscr{H}}^{s}$ and $s>\frac{d}{2}+1$. In the same way, $D^{2} G(u) \cdot\left(h_{1}, h_{2}\right)=\partial_{u}^{2} F(\cdot, u) \cdot\left(h_{1}, h_{2}\right)$ extends from $\tilde{\mathscr{H}}^{\sigma_{1}} \times \tilde{\mathscr{H}}^{\sigma_{2}}$ to $\tilde{\mathscr{H}}^{-\sigma_{3}}$ for $\sigma_{1}, \sigma_{2}, \sigma_{3}$ as in the statement of the definition, by Lemma A.3.
Definition 2.2.2. Let $\sigma \geq \sigma_{0}, X$ an open subset of $\mathscr{H}_{1}^{\sigma}, k \in \mathbb{Z}$. One denotes by $C^{\infty, k}(X ; \mathbb{R})$ the space of $C^{1}$ functions $\Phi: X \rightarrow \mathbb{R}$, such that for any $s \geq \sigma$ and $u \in X \cap \mathscr{H}_{1}^{s}$, we have $\nabla \Phi(u) \in \mathscr{H}_{1}^{s-k}$ and $u \rightarrow \nabla \Phi(u)$ belongs to $\Phi^{\infty, k}\left(X, \mathscr{H}_{1}^{\sigma-k}\right)$.

If $F: \mathbb{S}^{1} \times \mathbb{T}^{d} \times \mathbb{R}^{2} \rightarrow \mathbb{R}$ is a smooth function, with $F(t, x, 0) \equiv 0, \partial_{u} F(t, x, 0) \equiv 0, \partial_{u}^{2} F(t, x, 0) \equiv 0$, and if $\Phi(u)=\int_{\mathbb{S}^{1} \times \mathbb{T}^{d}} F(t, x, u(t, x)) d t d x$, then $\nabla \Phi(u)=\partial_{u} F(\cdot, u) \in \tilde{\mathscr{H}}^{s}$ if $u \in \tilde{\mathscr{H}}^{s}$ and $s>\frac{d}{2}+1$ (see Lemma A.1), and the example following Definition 2.2.1 shows that $\Phi \in C^{\infty, 0}\left(\tilde{\mathscr{H}}^{\sigma}, \mathbb{R}\right)$ for $\sigma \geq \sigma_{0}$.

Remark. In the sequel we shall have to consider elements $G(u, \omega, \epsilon), \Phi(u, \omega, \epsilon)$ of the preceding spaces depending on the real parameters $(\omega, \epsilon)$. We shall say that $G, \Phi$ are $C^{1}$ in $(\omega, \epsilon)$ if the conditions of Definition 2.2.1 (resp. Definition 2.2.2) are satisfied by $G, \partial_{\omega} G, \partial_{\epsilon} G$ (resp. $\Phi, \partial_{\omega} \Phi, \partial_{\epsilon} \Phi$ ).
Lemma 2.2.3. Let $\sigma \geq \sigma_{0}, k \in \mathbb{N}, X$ an open subset of $\mathscr{H}_{1}^{\sigma}, G \in \Phi^{\infty,-k}\left(X, \mathscr{H}_{2}^{\sigma+k}\right), Y$ an open subset of $\mathscr{H}_{2}^{\sigma+k}$ containing $G(X), \Phi \in C^{\infty, k}(Y, \mathbb{R})$. Then $\Phi \circ G \in C^{\infty, 0}(X, \mathbb{R})$.
Proof. The assumption on $G$ implies that for $v \in X \cap \mathscr{H}_{1}^{s}, s \geq \sigma$ and for $\sigma^{\prime}$ with $\left|\sigma^{\prime}\right| \leq s$,

$$
\begin{equation*}
D G(v) \in \mathscr{L}\left(\mathscr{H}_{1}^{\sigma^{\prime}}, \mathscr{H}_{2}^{\sigma^{\prime}+k}\right) \subset \mathscr{L}\left(\mathscr{H}_{1}^{\sigma^{\prime}}, \mathscr{H}_{2}^{\sigma^{\prime}}\right) \tag{2.2.1}
\end{equation*}
$$

Moreover, since $\nabla \Phi \in \Phi^{\infty, k}\left(Y, \mathscr{H}_{2}^{\sigma}\right)$, we have $G(v) \in Y \cap \mathscr{H}_{2}^{s+k}$ for $v \in X \cap \mathscr{H}_{1}^{s}$, so $\nabla \Phi(G(v)) \in \mathscr{H}_{2}^{s}$ and for any $\sigma^{\prime \prime}$ with $\left|\sigma^{\prime \prime}\right| \leq s+k,(D(\nabla \Phi))(G(v))$ is in $\mathscr{L}\left(\mathscr{H}_{2}^{\sigma^{\prime \prime}}, \mathscr{H}_{2}^{\sigma^{\prime \prime}-k}\right)$. In particular, for any $\sigma^{\prime}$ with $\left|\sigma^{\prime}\right| \leq s$,

$$
\begin{equation*}
D(\nabla \Phi)(G(v)) \in \mathscr{L}\left(\mathscr{H}_{2}^{\sigma^{\prime}+k}, \mathscr{H}_{2}^{\sigma^{\prime}}\right) \tag{2.2.2}
\end{equation*}
$$

We deduce from (2.2.1) that $\nabla(\Phi \circ G)(v)={ }^{t} D G(v) \cdot(\nabla \Phi)(G(v))$ belongs to $\mathscr{H}_{1}^{s}$ when $v \in X \cap \mathscr{H}_{1}^{s}$. Let us check that $\nabla(\Phi \circ G)$ belongs to $\Phi^{\infty, 0}\left(X, \mathscr{H}_{1}^{\sigma}\right)$. If $u \in X \cap \mathscr{H}_{1}^{s}(s \geq \sigma)$ and $h \in \mathscr{H}_{1}^{\sigma^{\prime}}$ with $\sigma^{\prime} \in[-s, s]$,
we write

$$
\begin{equation*}
D[\nabla(\Phi \circ G)(v)] \cdot h={ }^{t} D G(v) \cdot((D \nabla \Phi)(G(v)) \cdot D G(v) \cdot h)+\left(D\left({ }^{t} D G\right)(v) \cdot h\right) \cdot \nabla \Phi(G(v)) . \tag{2.2.3}
\end{equation*}
$$

By (2.2.1) and (2.2.2), the first term on the right belongs to $\mathscr{H}_{1}^{\sigma^{\prime}}$. To check that the last term in (2.2.3) belongs to the same space, we integrate it against $h^{\prime} \in \mathscr{H}_{1}^{-\sigma^{\prime}}$. We get

$$
\begin{equation*}
\int\left[\left(D\left({ }^{t} D G\right)(v) \cdot h\right) \cdot \nabla \Phi(G(v))\right] h^{\prime} d t d x=\int(\nabla \Phi)(G(v)) D^{2} G(v) \cdot\left(h, h^{\prime}\right) d t d x . \tag{2.2.4}
\end{equation*}
$$

By Definition 2.2.1,

$$
D^{2} G(v) \cdot\left(h, h^{\prime}\right) \in \mathscr{H}_{2}^{-\max \left(\sigma_{0}, \sigma^{\prime}\right)+k} \subset \mathscr{H}_{2}^{-\max \left(\sigma_{0}, \sigma^{\prime}\right)}
$$

Since $\nabla \Phi(G(v)) \in \mathscr{H}_{2}^{s} \subset \mathscr{H}_{2}^{\max \left(\sigma_{0}, \sigma^{\prime}\right)}$, this shows that the right side of (2.2.4) defines a continuous linear form in $h^{\prime} \in \mathscr{H}_{1}^{-\sigma^{\prime}}$.

We now study $D^{2}[\nabla(\Phi \circ G)(v)] \cdot\left(h_{1}, h_{2}\right)$, with $\left(h_{1}, h_{2}\right) \in \mathscr{H}_{1}^{\sigma_{1}} \times \mathscr{H}_{1}^{\sigma_{2}}$. To prove that

$$
D^{2}[\nabla(\Phi \circ G)(v)] \cdot\left(h_{1}, h_{2}\right) \in \mathscr{H}_{1}^{-\sigma_{3}}
$$

we compute, for $h_{3} \in \mathscr{H}_{1}^{\sigma_{3}}$,

$$
D^{2} \int \nabla(\Phi \circ G)(v) h_{3} d t d x=D^{2} \int[(\nabla \Phi)(G(v))]\left[D G(v) \cdot h_{3}\right] d t d x
$$

We get the following contributions (up to symmetries) for the action on $\left(h_{1}, h_{2}\right) \in \mathscr{H}_{1}^{\sigma_{1}} \times \mathscr{H}_{1}^{\sigma_{2}}$ :

$$
\begin{align*}
& \int[(\nabla \Phi)(G(v))]\left[D^{3} G(v) \cdot\left(h_{1}, h_{2}, h_{3}\right)\right] d t d x,  \tag{2.2.5a}\\
& \int\left[D((\nabla \Phi)(G(v))) \cdot h_{1}\right]\left[D^{2} G(v) \cdot\left(h_{2}, h_{3}\right)\right] d t d x,  \tag{2.2.5b}\\
& \int\left[(D \nabla \Phi)(G(v)) \cdot D^{2} G(v) \cdot\left(h_{1}, h_{2}\right)\right]\left[D G(v) \cdot h_{3}\right] d t d x,  \tag{2.2.5c}\\
& \int\left[\left(D^{2} \nabla \Phi\right)(G(v)) \cdot\left(D G(v) \cdot h_{1}, D G(v) \cdot h_{2}\right)\right]\left[D G(v) \cdot h_{3}\right] d t d x . \tag{2.2.5d}
\end{align*}
$$

In (2.2.5a), we may assume for instance $h_{1} \in \mathscr{H}_{1}^{\sigma^{\prime}}, h_{2} \in \mathscr{H}_{1}^{-\sigma^{\prime}}, h_{3} \in \mathscr{H}_{1}^{\max \left(\sigma_{0}, \sigma^{\prime}\right)}$. Since $u \rightarrow D^{2} G(u)$ is $C^{1}$ on $X \cap \mathscr{H}_{1}^{\max \left(\sigma_{0}, \sigma^{\prime}\right)}$ with values in $\mathscr{L}_{2}\left(\mathscr{H}_{1}^{\sigma^{\prime}} \times \mathscr{H}_{1}^{-\sigma^{\prime}} ; \mathscr{H}_{2}^{-\max \left(\sigma_{0}, \sigma^{\prime}\right)+k}\right)$, the second factor in the integrand belongs to $\mathscr{H}_{2}^{-\max \left(\sigma_{0}, \sigma^{\prime}\right)+k}$, so may be integrated against $\nabla \Phi(G(v)) \in \mathscr{H}_{2}^{s} \subset \mathscr{H}_{2}^{\max \left(\sigma_{0}, \sigma^{\prime}\right)}$ for $s \geq \sigma^{\prime} \geq 0$ and $s \geq \sigma$.

In $(2.2 .5 \mathrm{~b}), D^{2} G(v) \cdot\left(h_{2}, h_{3}\right) \in \mathscr{H}_{2}^{-\sigma_{1}+k}$. On the other hand $D((\nabla \Phi)(G(v))) \cdot h_{1}$ is in $\mathscr{H}_{2}^{\sigma_{1}}$ by (2.2.1) and (2.2.2), which allows one to integrate the product of the two factors.
In (2.2.5c), DG(v) $\cdot h_{3}$ lies in $\mathscr{H}_{2}^{\sigma_{3}+k}$. The other factor is given by the action of $(D \nabla \Phi)(G(v))$ on $D^{2} G(v) \cdot\left(h_{1}, h_{2}\right) \in \mathscr{H}_{2}^{-\sigma_{3}+k}$, whence again the wanted duality in the integral, using (2.2.2).

Finally, in $(2.2 .5 \mathrm{~d})$, we integrate $D G(v) \cdot h_{3} \in \mathscr{H}_{2}^{\sigma_{3}+k}$ against the action of $\left(D^{2} \nabla \Phi\right)(G(v))$ on a couple belonging to $\mathscr{H}_{2}^{\sigma_{1}+k} \times \mathscr{H}_{2}^{\sigma_{2}+k} \subset \mathscr{H}_{2}^{\sigma_{1}} \times \mathscr{H}_{2}^{\sigma_{2}}$. Since this vector is in $\mathscr{H}_{2}^{-\sigma_{3}-k}$ by definition of $C^{\infty, k}(Y, \mathbb{R})$, we get the conclusion.

Let us write an equivalent form of (1.1.3) using the classes of functions above. Since the Hamiltonian $F$ in (1.1.2) is real-valued, we may write (1.1.3) as a $2 \times 2$ system

$$
\begin{align*}
\left(\omega D_{t}-\Delta+\mu\right) u & =\epsilon f(t, x)+\epsilon \frac{\partial F}{\partial \bar{u}}(t, x, u, \bar{u}, \epsilon)  \tag{2.2.6}\\
\left(-\omega D_{t}-\Delta+\mu\right) \bar{u} & =\epsilon \bar{f}(t, x)+\epsilon \frac{\partial F}{\partial u}(t, x, u, \bar{u}, \epsilon)
\end{align*}
$$

We identify $u=v_{1}+i v_{2}$ with $v=\left[\begin{array}{l}v_{1} \\ v_{2}\end{array}\right]$ and $f=f_{1}+i f_{2}$ with $f=\left[\begin{array}{l}f_{1} \\ f_{2}\end{array}\right]$. If we set

$$
\nabla F(v)=\left[\begin{array}{l}
\partial F / \partial v_{1} \\
\partial F / \partial v_{2}
\end{array}\right]
$$

and

$$
L_{\omega}=\left[\begin{array}{cc}
\Delta-\mu & -\omega \partial_{t}  \tag{2.2.7}\\
\omega \partial_{t} & \Delta-\mu
\end{array}\right]
$$

Equation (2.2.6) is equivalent to

$$
\begin{equation*}
L_{\omega} v=-\epsilon f-\epsilon \nabla_{v} F(t, x, v) \tag{2.2.8}
\end{equation*}
$$

Define, for $v \in \widetilde{\mathscr{H}}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right)$,

$$
\begin{equation*}
\Phi_{1}(v, f, \omega, \epsilon)=\frac{1}{2} \int_{\mathbb{S}^{1} \times \mathbb{T}^{d}}\left(L_{\omega} v\right) v d t d x+\epsilon \int_{\mathbb{S}^{1} \times \mathbb{T}^{d}} f(t, x) v(t, x) d t d x \tag{2.2.9}
\end{equation*}
$$

and

$$
\begin{equation*}
\Phi_{2}(v, \epsilon)=\int_{\mathbb{S}^{1} \times \mathbb{T}^{d}} F(t, x, v(t, x), \epsilon) d t d x \tag{2.2.10}
\end{equation*}
$$

Then $\nabla \Phi_{1}(v)=L_{\omega} v+\epsilon f$, so $\Phi_{1} \in C^{\infty, 2}\left(\tilde{\mathscr{H}}^{\sigma} \times \tilde{\mathscr{H}}^{\sigma}, \mathbb{R}\right)$ if $\sigma \geq \sigma_{0}$, since, by the definition of $\tilde{\mathscr{H}}^{\sigma}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right)$, $L_{\omega}$ is bounded from $\tilde{\mathscr{H}}^{\sigma}$ to $\tilde{\mathscr{H}}^{\sigma-2}$. By the statement following Definition 2.2.2, $\Phi_{2} \in C^{\infty, 0}\left(\tilde{\mathscr{H}}^{\sigma}, \mathbb{R}\right)$ ( $\sigma \geq \sigma_{0}$ ). Moreover (2.2.8) may be written

$$
\begin{equation*}
\nabla_{v}\left[\Phi_{1}(v, f, \omega, \epsilon)+\epsilon \Phi_{2}(v, \epsilon)\right]=0 \tag{2.2.11}
\end{equation*}
$$

Using the notation introduced at the bottom of page 646, we decompose any $v \in \tilde{\mathscr{H}}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right)$ as $v=v^{\prime}+v^{\prime \prime}$ on the decomposition

$$
\tilde{\mathscr{H}}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right)=\mathscr{H}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right) \oplus \mathscr{F}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right)
$$

We denote for $q>0$ by $B_{q}\left(\tilde{\mathscr{H}}^{s}\right), B_{q}\left(\mathscr{H}^{s}\right), B_{q}\left(\mathscr{F}^{s}\right)$ the ball of center 0 and radius $q$ in these spaces. By (1.2.6), if $v \in \mathscr{F}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right),(j, n) \in \mathbb{Z} \times \Omega_{\alpha} \subset \mathbb{Z} \times \mathbb{Z}^{d}$ and $\hat{v}(j, n) \neq 0$, then $|j|>K_{0}\langle n(\alpha)\rangle^{2}$ or $|j|<K_{0}^{-1}\langle n(\alpha)\rangle^{2}$. Moreover, since $\mu \in \mathbb{R}-\mathbb{Z}_{-},\left||n|^{2}+\mu\right| \geq c(\mu)\langle n(\alpha)\rangle^{2}$ when $n \in \Omega_{\alpha}$, for some constant $c(\mu)>0$. If we fix $K_{0}$ large enough, and use that $\omega$ stays in [1, 2], we conclude that the eigenvalues of $L_{\omega}$ satisfy the bounds

$$
\left|\omega j+|n|^{2}+\mu\right| \geq c\left(|j|+\langle n(\alpha)\rangle^{2}\right) \quad \text { for } j \in \mathbb{Z}, n \in \Omega_{\alpha}, \alpha \in \mathscr{A}
$$

This shows that the restriction of $L_{\omega}$ to $\mathscr{F}^{s+2}$ is an invertible operator from $\mathscr{F}^{s+2}$ to $\mathscr{F}^{s}$ (uniformly in $\omega \in[1,2]$ ).

Let us reduce (2.2.11) to an equation on the space $\mathscr{H}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right)$.
Proposition 2.2.4. Let $\sigma \geq \sigma_{0}, q>0, f^{\prime} \in B_{q}\left(\mathscr{H}^{\sigma}\right)$. There are $\left.\left.\gamma_{0} \in\right] 0,1\right]$,

- an element $\left(v^{\prime}, f^{\prime \prime}\right) \rightarrow \psi_{2}\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right)$ of $C^{\infty, 0}\left(W_{q} ; \mathbb{R}\right)$, where $W_{q}=B_{q}\left(\mathscr{H}^{\sigma}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right)\right) \times$ $B_{q}\left(\mathscr{F}^{\sigma}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right)\right)$, with $C^{1}$ dependence in $(\omega, \epsilon) \in[1,2] \times\left[0, \gamma_{0}\right]$, and
- an element $\left(v^{\prime}, f^{\prime \prime}\right) \rightarrow G\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right)$ of $\Phi^{\infty,-2}\left(W_{q}, \mathscr{F}^{\sigma+2}\right)$, with $C^{1}$ dependence in $(\omega, \epsilon)$,
such that, for any given subset $A \subset[1,2] \times\left[0, \gamma_{0}\right]$, the following two conditions are equivalent:
(i) The function $v=\left(v^{\prime}, G\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right)\right)$ satisfies for any $(\omega, \epsilon) \in A$

$$
\begin{equation*}
L_{\omega} v+\epsilon f+\epsilon \nabla_{v} \Phi_{2}(v, \epsilon)=0 \tag{2.2.12}
\end{equation*}
$$

where $f=f^{\prime}+f^{\prime \prime}$.
(ii) The function $v^{\prime}$ satisfies for any $(\omega, \epsilon) \in A$

$$
\begin{equation*}
L_{\omega} v^{\prime}+\epsilon f^{\prime}+\epsilon \nabla_{v^{\prime}} \psi_{2}\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right)=0 \tag{2.2.13}
\end{equation*}
$$

Proof. Write (2.2.12) as

$$
\begin{align*}
L_{\omega} v^{\prime}+\epsilon f^{\prime}+\epsilon \nabla_{v^{\prime}} \Phi_{2}\left(v^{\prime}, v^{\prime \prime}, \epsilon\right) & =0  \tag{2.2.14a}\\
L_{\omega} v^{\prime \prime}+\epsilon f^{\prime \prime}+\epsilon \nabla_{v^{\prime \prime}} \Phi_{2}\left(v^{\prime}, v^{\prime \prime}, \epsilon\right) & =0 \tag{2.2.14b}
\end{align*}
$$

We are looking for a solution of the second equation under the form $v^{\prime \prime}=-\epsilon L_{\omega}^{-1} f^{\prime \prime}+\epsilon w^{\prime \prime}$. The new unknown $w^{\prime \prime}$ satisfies

$$
\begin{equation*}
w^{\prime \prime}=-L_{\omega}^{-1} \nabla_{v^{\prime \prime}} \Phi_{2}\left(v^{\prime},-\epsilon L_{\omega}^{-1} f^{\prime \prime}+\epsilon w^{\prime \prime}, \epsilon\right) \tag{2.2.15}
\end{equation*}
$$

Let $q_{0}>0$ be such that $\left\|L_{\omega}^{-1} \nabla_{v^{\prime \prime}} \Phi_{2}\left(v^{\prime}, h, \epsilon\right)\right\|_{\mathscr{F} \sigma+2} \leq q_{0} / 2$ for any $\left(v^{\prime}, h\right) \in B_{q}\left(\mathscr{H}^{\sigma}\right) \times B_{q}\left(\mathscr{F}^{\sigma}\right)$, any $\epsilon \in[0,1]$, any $\omega \in[1,2]$. The fixed point theorem with parameters shows that there is $\left.\left.\gamma_{0} \in\right] 0,1\right]$ such that for any $\left(v^{\prime}, f^{\prime \prime}\right) \in W_{q}$, any $\epsilon \in\left[0, \gamma_{0}\right]$, Equation (2.2.15) has a unique solution $w^{\prime \prime} \in B_{q_{0}}\left(\mathscr{F}^{\sigma+2}\right)$. We denote this solution by $G\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right)$. This is a smooth function of $\left(v^{\prime}, f^{\prime \prime}\right) \in W_{q}$, with $C^{1}$ dependence in $(\omega, \epsilon)$. If moreover $\left(v^{\prime}, f^{\prime \prime}\right) \in \tilde{\mathscr{H}}^{s}$ for some $s \geq \sigma$, it follows from (2.2.15) that $w^{\prime \prime} \in \mathscr{F}^{s+2}$ (using that $L_{\omega}^{-1}$ gains two derivatives in the $\mathscr{F}^{s}$ scale $)$. Let us show that $G$ belongs to $\Phi^{\infty,-2}\left(W_{q}, \mathscr{F}^{\sigma+2}\right)$. By the definition of $G$

$$
\begin{align*}
D_{v^{\prime}} G\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right) & =-L_{\omega}^{-1}\left(\operatorname{Id}-\epsilon M^{\prime \prime}\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right) L_{\omega}^{-1}\right)^{-1} M^{\prime}\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right)  \tag{2.2.16}\\
D_{f^{\prime \prime}} G\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right) & =\epsilon L_{\omega}^{-1}\left(\operatorname{Id}-\epsilon M^{\prime \prime}\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right) L_{\omega}^{-1}\right)^{-1} M^{\prime \prime}\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right) L_{\omega}^{-1}
\end{align*}
$$

with

$$
\begin{align*}
M^{\prime}\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right) & =\left(D_{v^{\prime}} \nabla_{v^{\prime \prime}} \Phi_{2}\right)\left(v^{\prime},-\epsilon L_{\omega}^{-1} f^{\prime \prime}+\epsilon G, \epsilon\right) \\
M^{\prime \prime}\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right) & =-\left(D_{v^{\prime \prime}} \nabla_{v^{\prime \prime}} \Phi_{2}\right)\left(v^{\prime},-\epsilon L_{\omega}^{-1} f^{\prime \prime}+\epsilon G, \epsilon\right) \tag{2.2.17}
\end{align*}
$$

Since $\Phi_{2} \in C^{\infty, 0}\left(W_{q}, \mathbb{R}\right)$, when $\left(v^{\prime}, f^{\prime \prime}\right) \in W_{q} \cap \tilde{\mathscr{H}}^{s}$ for some $s \geq \sigma$, one can extend $M^{\prime \prime}\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right)$ into an element of $\mathscr{L}\left(\mathscr{F} \sigma^{\prime}, \mathscr{F} \sigma^{\prime}\right)$ for any $\sigma^{\prime} \in[-s, s]$; similarly, $M^{\prime}\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right)$ extends as an element
of $\mathscr{L}\left(\mathscr{H} \sigma^{\sigma^{\prime}}, \mathscr{F} \sigma^{\prime}\right)$. We choose $\gamma_{0}$ small enough that for $\left.\epsilon \in\left[0, \gamma_{0}\right], \epsilon\left\|M^{\prime \prime}\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right) L_{\omega}^{-1}\right\|_{\mathscr{L}(\mathscr{F} \sigma}, \mathscr{F} \sigma\right)$ is smaller than $\frac{1}{2}$. Let us check that $G$ satisfies the first condition in Definition 2.1.1. We may write the first equation in (2.2.16) as

$$
\begin{align*}
& D_{v^{\prime}} G\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right) \\
& \quad=-\sum_{k=0}^{2 N-1} L_{\omega}^{-1}\left(\epsilon M^{\prime \prime} L_{\omega}^{-1}\right)^{k} M^{\prime}-L_{\omega}^{-1}\left(\epsilon M^{\prime \prime} L_{\omega}^{-1}\right)^{N}\left(\operatorname{Id}-\epsilon M^{\prime \prime} L_{\omega}^{-1}\right)^{-1}\left(\epsilon M^{\prime \prime} L_{\omega}^{-1}\right)^{N} M^{\prime} \tag{2.2.18}
\end{align*}
$$

and a similar formula holds for $D_{f^{\prime \prime}} G$. If $N$ is chosen large enough relatively to $s$, and $\sigma^{\prime} \in[-s, s]$, $\left(\epsilon M^{\prime \prime} L_{\omega}^{-1}\right)^{N} M^{\prime}$ sends $\mathscr{H}^{\sigma^{\prime}}$ to $\mathscr{F} \sigma$, over which $\left(\operatorname{Id}-\epsilon M^{\prime \prime} L_{\omega}^{-1}\right)^{-1}$ is bounded. Consequently, the last contribution in (2.2.18) is in $\mathscr{F}^{s+2} \subset \mathscr{F}^{\sigma^{\prime}+2}$. The sum on the right side being bounded from $\mathscr{H}^{\sigma^{\prime}}$ to $\mathscr{F} \sigma^{\sigma^{\prime}+2}$ for any $\sigma^{\prime} \in[-s, s]$, we get the same property for $D_{v^{\prime}} G$. We argue in the same way for $D_{f^{\prime \prime}} G$. To check the second condition in Definition 2.1.1, we compute from (2.2.16), for $\left(h_{1}, h_{2}\right) \in \mathscr{H}^{\sigma_{1}} \times \mathscr{H}^{\sigma_{2}}$

$$
\begin{aligned}
D_{v^{\prime}}^{2} G\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right) \cdot\left(h_{1}, h_{2}\right)= & -L_{\omega}^{-1}\left(\operatorname{Id}-\epsilon M^{\prime \prime} L_{\omega}^{-1}\right)^{-1}\left[\left(D_{v^{\prime}} M^{\prime} \cdot h_{1}\right) \cdot h_{2}\right] \\
& -L_{\omega}^{-1}\left(\operatorname{Id}-\epsilon M^{\prime \prime} L_{\omega}^{-1}\right)^{-1}\left(\epsilon D_{v^{\prime}} M^{\prime \prime} L_{\omega}^{-1} \cdot h_{1}\right)\left(\operatorname{Id}-\epsilon M^{\prime \prime} L_{\omega}^{-1}\right)^{-1} M^{\prime} \cdot h_{2}
\end{aligned}
$$

If $\left\{\sigma_{1}, \sigma_{2}, \sigma_{3}\right\}=\left\{\sigma^{\prime},-\sigma^{\prime}, \max \left(\sigma_{0}, \sigma^{\prime}\right)\right\}$, the assumption on $\Phi_{2}$ implies that $D_{v^{\prime}} M^{\prime}$ sends $\mathscr{H}^{\sigma_{1}} \times \mathscr{H}^{\sigma_{2}}$ to $\mathscr{F}^{-\sigma_{3}}$, and $D_{v^{\prime}} M^{\prime \prime}$ sends $\mathscr{H}^{\sigma_{1}} \times \mathscr{F}^{\sigma_{2}}$ to $\mathscr{F}^{-\sigma_{3}}$. Using expansions as in (2.2.18), we conclude that if $\left(h_{1}, h_{2}\right) \in \mathscr{H}^{\sigma_{1}} \times \mathscr{H}^{\sigma_{2}}, D_{v^{\prime}}^{2} G\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right) \cdot\left(h_{1}, h_{2}\right) \in \mathscr{F}-\sigma_{3}+2$. One studies $D_{v^{\prime}} D_{f^{\prime \prime}} G$ and $D_{f^{\prime \prime}}^{2} G$ in the same way. It is clear that $D G, D^{2} G$ are smooth in $\left(v^{\prime}, f^{\prime \prime}\right) \in W_{q} \cap \widetilde{\mathscr{H}}^{s}$ and have a $C^{1}$ dependence in $(\omega, \epsilon)$; hence $G \in \Phi^{\infty,-2}\left(W_{q}, \mathscr{F}^{\sigma+2}\right)$.

Let us obtain the equivalent form (2.2.13) of (2.2.12) or (2.2.11). By (2.2.9), (2.2.10)

$$
\begin{aligned}
& \Phi_{1}\left(v^{\prime}, v^{\prime \prime}, \omega, \epsilon\right)+\epsilon \Phi_{2}\left(v^{\prime}, v^{\prime \prime}, \epsilon\right)=\frac{1}{2} \int\left(L_{\omega} v^{\prime}\right) v^{\prime} d t d x+\epsilon \int f^{\prime} v^{\prime} d t d x \\
&+\frac{1}{2} \int\left(L_{\omega} v^{\prime \prime}\right) v^{\prime \prime} d t d x+\epsilon \int f^{\prime \prime} v^{\prime \prime} d t d x+\epsilon \Phi_{2}\left(v^{\prime}, v^{\prime \prime}, \epsilon\right)
\end{aligned}
$$

We plug into this expression the solution $v^{\prime \prime}=-\epsilon L_{\omega}^{-1} f^{\prime \prime}+\epsilon G\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right)$ of (2.2.14b). We get after simplification the function
$\Psi\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right)=\frac{1}{2} \int\left(L_{\omega} v^{\prime}\right) v^{\prime} d t d x+\epsilon \int f^{\prime} v^{\prime} d t d x-\frac{\epsilon^{2}}{2} \int\left(L_{\omega}^{-1} f^{\prime \prime}\right) f^{\prime \prime} d t d x+\epsilon \psi_{2}\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right)$,
where

$$
\begin{equation*}
\psi_{2}\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right)=\frac{\epsilon}{2} \int G\left(L_{\omega} G\right) d t d x+\Phi_{2}\left(v^{\prime},-\epsilon L_{\omega}^{-1} f^{\prime \prime}+\epsilon G, \epsilon\right) \tag{2.2.19}
\end{equation*}
$$

The integral in (2.2.19) is the composition of the function defined on $\mathscr{F}^{\sigma}$ by $w^{\prime \prime} \rightarrow \int w^{\prime \prime}\left(L_{\omega} w^{\prime \prime}\right) d t d x$, which is an element of $C^{\infty, 2}\left(\mathscr{F}^{\sigma}, \mathbb{R}\right)$, with the map

$$
\begin{aligned}
\left(v^{\prime}, f^{\prime \prime}\right) & \rightarrow G\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right), \\
\tilde{\mathscr{H}}^{\sigma} & \rightarrow \mathscr{F}^{\sigma+2}
\end{aligned}
$$

which is an element of $\Phi^{\infty,-2}\left(W_{q}, \mathscr{F}^{\sigma+2}\right)$. By Lemma 2.2.3, we conclude that $\psi_{2} \in C^{\infty, 0}\left(W_{q}, \mathbb{R}\right)$.

Since $G$ is defined as the critical point (up to an affine change of variables) of the map

$$
v^{\prime \prime} \rightarrow\left(\Phi_{1}+\epsilon \Phi_{2}\right)\left(v^{\prime}, v^{\prime \prime}, \omega, \epsilon\right)
$$

and since $\Psi$ is the corresponding critical value, $v^{\prime}$ solves (2.2.14a) if and only of $\nabla_{v^{\prime}} \Psi\left(v^{\prime}, f^{\prime \prime}, \omega, \epsilon\right)=0$. This gives (2.2.13).

We finish this subsection with a lemma that will be useful in the sequel. Let $X$ be an open subset of $\mathscr{H}^{\sigma_{0}}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right), \psi$ an element of $C^{\infty, 0}(X ; \mathbb{R})$. For $v \in X \cap \mathscr{H}^{+\infty}, w_{1}, w_{2} \in \mathscr{H}^{+\infty}$, we set

$$
\begin{equation*}
L\left(v ; w_{1}, w_{2}\right)=D^{2} \psi(v) \cdot\left(w_{1}, w_{2}\right) \tag{2.2.20}
\end{equation*}
$$

This is a continuous bilinear form in $\left(w_{1}, w_{2}\right) \in \mathscr{H}^{0} \times \mathscr{H}^{0}$, by the definition of $C^{\infty, 0}(X ; \mathbb{R})$. By the Riesz theorem, we write it

$$
L\left(v ; w_{1}, w_{2}\right)=\int_{\mathbb{S}^{1} \times \mathbb{T}^{d}}\left(W(v) w_{1}\right) w_{2} d t d x
$$

for some symmetric $\mathscr{H}^{0}$-bounded operator $W(v)$. Since Definition 2.2.2 implies that $v \rightarrow D^{2} \psi(v)$ is a smooth map on $X$ with values in the space of continuous bilinear forms on $\mathscr{H}^{0} \times \mathscr{H}^{0}$, we know that $v \rightarrow W(v)$ is smooth with values in $\mathscr{L}\left(\mathscr{H}^{0}, \mathscr{H}^{0}\right)$. Thus we may write, for $j=1, \ldots, d$,

$$
\begin{align*}
L\left(v ; \partial_{x_{j}} w_{1}, w_{2}\right)+L\left(v ; w_{1}, \partial_{x_{j}} w_{2}\right) & =-\int_{\mathbb{S}^{1} \times \mathbb{T}^{d}}\left(\left(\partial_{x_{j}} W(v)\right) w_{1}\right) w_{2} d t d x \\
& =-\left(\partial_{v} L\right)\left(v ; w_{1}, w_{2}\right) \cdot\left(\partial_{x_{j}} v\right) \tag{2.2.21}
\end{align*}
$$

for any $v \in X \cap \mathscr{H}^{+\infty}$ and $w_{1}, w_{2} \in \mathscr{H}^{+\infty}$.
We denote by $\mathbb{C}\left[X_{\alpha} ; \alpha \in \mathbb{N}^{d}\right]$ the space of polynomials in indeterminates $X_{\alpha}$, indexed by elements of $\mathbb{N}^{d}$. If $X_{\alpha_{1}}^{k_{1}} \cdots X_{\alpha_{\ell}}^{k_{\ell}}$ is a monomial, its weight will be defined as $k_{1}\left|\alpha_{1}\right|+\cdots+k_{\ell}\left|\alpha_{\ell}\right|$. The weight of any polynomial is then defined in the natural way.

Lemma 2.2.5. For any $N \in \mathbb{N}$ and $\ell \in \mathbb{N}$, there is a polynomial $Q_{N}^{\ell} \in \mathbb{C}\left[X_{\alpha} ; \alpha \in \mathbb{N}^{d}\right]$, of weight less or equal to $N$, and for any $q>0$ a constant $C>0$ such that, for any $v \in B_{q}\left(\mathscr{H}^{\sigma_{0}}\right) \cap \mathscr{H}^{+\infty} \cap X$, any $h_{1}, \ldots, h_{\ell}$ in $\mathcal{H}^{+\infty}$, any $n, n^{\prime} \in \mathbb{Z}^{d}$

$$
\begin{align*}
& \left.\left\|\Pi_{n} \partial_{v}^{\ell} W(v) \cdot\left(h_{1}, \ldots, h_{\ell}\right) \Pi_{n^{\prime}}\right\|_{\mathscr{L}(\mathscr{H} 0}\right) \\
& \leq C\left\langle n-n^{\prime}\right\rangle^{-N} \sum_{N_{0}+\cdots+N_{\ell}=N} Q_{N_{0}}^{\ell}\left(\left(\left\|\partial^{\alpha} v\right\|_{\mathscr{H}} \sigma_{0}\right)_{\alpha}\right) \prod_{\ell^{\prime}=1}^{\ell}\left\|h_{\ell^{\prime}}\right\|_{\mathscr{H}^{\sigma} \sigma_{0}+N_{\ell^{\prime}}} . \tag{2.2.22}
\end{align*}
$$

Proof. Since ${ }^{t} \Pi_{n}=\Pi_{-n}$, we may write, for any $w_{1}, w_{2} \in \mathscr{H}^{+\infty}$,

$$
\begin{aligned}
\left(n_{j}-n_{j}^{\prime}\right) \int\left(\Pi_{n} W(v) \Pi_{n^{\prime}} w_{1}\right) w_{2} d t d x & =\left(n_{j}-n_{j}^{\prime}\right) L\left(v ; \Pi_{n^{\prime}} w_{1}, \Pi_{-n} w_{2}\right) \\
& =i\left[L\left(v ; \partial_{x_{j}} \Pi_{n^{\prime}} w_{1}, \Pi_{-n} w_{2}\right)+L\left(v ; \Pi_{n^{\prime}} w_{1}, \partial_{x_{j}} \Pi_{-n} w_{2}\right)\right] \\
& =-i\left(\partial_{v} L\right)\left(v ; \Pi_{n^{\prime}} w_{1}, \Pi_{-n} w_{2}\right) \cdot\left(\partial_{x_{j}} v\right)
\end{aligned}
$$

by (2.2.21). Iterating the computation, we get for

$$
\left\langle n-n^{\prime}\right\rangle^{N}\left|\int\left(\Pi_{n} W(v) \Pi_{n^{\prime}} w_{1}\right) w_{2} d t d x\right|
$$

an estimate in terms of quantities

$$
\left|\left(\partial_{v}^{p} L\right)\left(v ; \Pi_{n^{\prime}} w_{1}, \Pi_{-n} w_{2}\right) \cdot\left(\partial^{\alpha_{1}} v, \ldots, \partial^{\alpha_{p}} v\right)\right|
$$

with $\left|\alpha_{1}\right|+\cdots+\left|\alpha_{p}\right| \leq N$. By the properties of $L$, this is bounded from above by

$$
C\left\|\Pi_{n^{\prime}} w_{1}\right\|_{L^{2}}\left\|\Pi_{-n} w_{2}\right\|_{L^{2}} \prod_{p^{\prime}=1}^{p}\left\|\partial^{\alpha_{p^{\prime}}} v\right\|_{\mathscr{H}} \sigma_{0}
$$

when $v$ stays in a fixed $\mathscr{H}^{\sigma_{0}}$-ball. This implies (2.2.22) for $\ell=0$. The proof for general $\ell$ is similar, up to notation.
2.3. Reduction to a paradifferential equation. We want to construct, under the conditions of the statement of Theorem 1.1.1, periodic solutions to (2.2.6). We have rewritten this equation under the real form (2.2.8) (or (2.2.11)). By Proposition 2.2.4, if we find a periodic solution $v^{\prime}$ for (2.2.13), we get a periodic solution $v$ for (2.2.12), which is a rewriting of (2.2.11). We are thus reduced to finding a solution $v^{\prime} \in \mathscr{H}^{\sigma}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right)$ to (2.2.13). Since the force term $f=f^{\prime}+f^{\prime \prime}$ will be fixed, we no longer write the $f^{\prime \prime}$ dependence in the function $\psi_{2}$ defined in Proposition 2.2.4. Moreover, since, in the rest of the paper, we will study only the equivalent formulation (2.2.13) of our initial problem, we drop the primes; that is, we study

$$
\begin{equation*}
L_{\omega} v+\epsilon f+\epsilon \nabla_{v} \psi_{2}(v, \omega, \epsilon)=0 \tag{2.3.1}
\end{equation*}
$$

where $v \in B_{q}\left(\mathscr{H}^{\sigma}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right)\right), f \in \mathscr{H}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right), \psi_{2}$ is in $C^{\infty, 0}\left(B_{q}\left(\mathscr{H}^{\sigma}\right), \mathbb{R}\right)$ for some $\sigma \in\left[\sigma_{0}, s\right]$, $q>0$ and for $\epsilon \in\left[0, \gamma_{0}\right]$, with $\left.\left.\gamma_{0} \in\right] 0,1\right]$ small enough. We shall use the equivalent norms (1.2.7) and (1.2.8) on the spaces we consider.

Our objective in this subsection is to rewrite the nonlinearity in (2.3.1) using paradifferential operators.
Proposition 2.3.1. Let $q>0, \sigma \geq \sigma_{0}+d+1$ be given. Set

$$
\begin{equation*}
r=\sigma-\sigma_{0}-d-1 \tag{2.3.2}
\end{equation*}
$$

There is a symmetric element $\tilde{V} \in \Psi_{\mathbb{R}}^{0}(0, \sigma, q) \otimes \mathcal{M}_{2}(\mathbb{R})$ and an element $\widetilde{R} \in \mathscr{R}_{0, \mathbb{R}}^{r}(0, \sigma, q) \otimes \mathcal{M}_{2}(\mathbb{R})$, with $C^{1}$ dependence in $(\omega, \epsilon)$, such that, for any $v \in B_{q}\left(\mathscr{H}^{\sigma}\right), \epsilon \in\left[0, \gamma_{0}\right]$, and $\omega \in[1,2]$,

$$
\begin{equation*}
\nabla_{v} \psi_{2}(v, \omega, \epsilon)=\tilde{V}(v, \omega, \epsilon) v+\widetilde{R}(v, \omega, \epsilon) v \tag{2.3.3}
\end{equation*}
$$

Let us comment about the interest of this decomposition of $\nabla_{v} \psi_{2}$. It allows us to express the nonlinearity in $(2.3 .1)$ as the sum of a remainder and of the action of the paradifferential potential $\tilde{V}(v, \omega, \epsilon)$ on $v$. In that way, the main contribution to the nonlinearity is expressed in terms of a class of operators enjoying a nice calculus. This will be exploited below to perform a block diagonalization.

We introduce some notation for the proof. For $p \in \mathbb{N}, v \in \mathscr{H}^{0}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right)$, we set

$$
\begin{array}{lll}
\Delta_{0} v=\Pi_{0} v, & \Delta_{p} v=\sum_{\substack{n \in \mathbb{Z}^{d} \\
2^{p-1} \leq|n|<2^{p}}} \Pi_{n} v & \text { for } p \geq 1 \\
S_{0} v=0, & S_{p} v=\sum_{p^{\prime}=0}^{p-1} \Delta_{p^{\prime}} v=\sum_{\substack{n \in \mathbb{Z}^{d} \\
|n|<2^{p-1}}} \Pi_{n} v \quad \text { for } p \geq 1 \tag{2.3.4}
\end{array}
$$

We also consider the frequency cut-offs defined for $n, n^{\prime} \in \mathbb{Z}^{d}$ by

$$
\begin{equation*}
S\left(n, n^{\prime}\right)=\sum_{\left|n^{\prime \prime}\right| \leq 2\left(1+\min \left(|n|,\left|n^{\prime}\right|\right)\right)} \Pi_{n^{\prime \prime}} \tag{2.3.5}
\end{equation*}
$$

Lemma 2.3.2. Let $\sigma \geq \sigma_{0}+d+1, q>0$. There is a map $(v, \omega, \epsilon) \rightarrow W(v, \omega, \epsilon)$ defined for $v \in B_{q}\left(\mathscr{H}^{\sigma}\right)$, $\epsilon \in\left[0, \gamma_{0}\right], \omega \in[1,2]$, with values in the space of bounded symmetric operators on $\mathscr{H}^{0}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right)$, which is $C^{\infty}$ in $v$ and has $C^{1}$ dependence in ( $\omega, \epsilon$ ), such that for any $(v, \omega, \epsilon$ )

$$
\begin{equation*}
\psi_{2}(v, \omega, \epsilon)=\int_{\mathbb{S}^{1} \times \mathbb{T}^{d}}[W(v, \omega, \epsilon) v] v d t d x \tag{2.3.6}
\end{equation*}
$$

and such that the following estimate holds: For $(\ell, N) \in \mathbb{N} \times \mathbb{N}$, there are polynomials $Q_{N}^{\ell} \in \mathbb{C}\left[X_{\alpha} ; \alpha \in \mathbb{N}\right]$, of weight at most $N$, and there is for any $M \in \mathbb{N}$ and $\ell \in \mathbb{N}$ a constant $C$, depending only on $\ell, q, M$, such that for any $v \in B_{q}\left(\mathscr{H}^{\sigma}\right)$, any $\epsilon \in\left[0, \gamma_{0}\right]$, any $\omega \in[1,2]$, any $\left(a_{0}, a_{1}\right) \in \mathbb{N}^{2}$ with $a_{0}+a_{1} \leq 1$, any $\left(h_{1}, \ldots, h_{\ell}\right) \in\left(\mathscr{H}^{\sigma}\right)^{\ell}$, and any $n, n^{\prime} \in \mathbb{Z}^{d}$,

$$
\left\|\Pi_{n} \partial_{\omega}^{a_{0}} \partial_{\epsilon}^{a_{1}} D_{v}^{\ell} W(v, \omega, \epsilon) \cdot\left(h_{1}, \ldots, h_{\ell}\right) \Pi_{n^{\prime}}\right\|_{\mathscr{L}(\mathscr{H})}
$$

$$
\begin{equation*}
\leq C\left\langle n-n^{\prime}\right\rangle^{-M} \sum_{N_{0}+\cdots+N_{\ell}=M} Q_{N_{0}}^{\ell}\left(\left(\left\|\partial^{\alpha} S\left(n, n^{\prime}\right) v\right\|_{\mathscr{H}^{\sigma_{0}}}\right)_{\alpha}\right) \prod_{\ell^{\prime}=1}^{\ell}\left\|S\left(n, n^{\prime}\right) h_{\ell^{\prime}}\right\|_{\mathscr{H}^{\sigma_{0}+N_{\ell^{\prime}}}} . \tag{2.3.7}
\end{equation*}
$$

Proof. We do not write $\omega, \epsilon$, which play the role of parameters. Since $\psi_{2}$ vanishes at order 3 at $v=0$, and $S_{p} v \rightarrow v$ in $\mathscr{H}^{\sigma}$ when $p \rightarrow+\infty$, we write

$$
\psi_{2}(v)=\sum_{p_{1}=0}^{+\infty}\left(\psi_{2}\left(S_{p_{1}+1} v\right)-\psi_{2}\left(S_{p_{1}} v\right)\right)=\sum_{p_{1}=0}^{+\infty} \int_{0}^{1}\left(\partial \psi_{2}\right)\left(S_{p_{1}} v+\tau_{1} \Delta_{p_{1}} v\right) d \tau_{1} \cdot \Delta_{p_{1}} v .
$$

Repeating the process, we get

$$
\psi_{2}(v)=\sum_{p_{1}=0}^{+\infty} \sum_{p_{2}=0}^{+\infty} \int_{0}^{1} \int_{0}^{1}\left(\partial^{2} \psi_{2}\right)\left(\Omega_{p_{1}, p_{2}}\left(\tau_{1}, \tau_{2}\right) v\right) d \tau_{2} \cdot\left(\Delta_{p_{2}}\left(S_{p_{1}}+\tau_{1} \Delta_{p_{1}}\right) v, \Delta_{p_{1}} v\right) d \tau_{1},
$$

where $\Omega_{p_{1}, p_{2}}\left(\tau_{1}, \tau_{2}\right)=\prod_{\ell=1}^{2}\left(S_{p_{\ell}}+\tau_{\ell} \Delta_{p_{\ell}}\right)$. By the discussion before Lemma 2.2.5, there is a symmetric operator $\widetilde{W}(v)$ satisfying (2.2.22), such that

$$
\partial^{2} \psi_{2}(v) \cdot\left(w_{1}, w_{2}\right)=\int\left[\tilde{W}(v) w_{1}\right] w_{2} d t d x .
$$

We set

$$
\begin{align*}
W(v)=\frac{1}{2} \sum_{p_{1}} \sum_{p_{2}} & \int_{0}^{1} \int_{0}^{1} \Delta_{p_{1}}\left[\tilde{W}\left(\Omega_{p_{1}, p_{2}}\left(\tau_{1}, \tau_{2}\right) v\right) \Delta_{p_{2}}\left(S_{p_{1}}+\tau_{1} \Delta_{p_{1}}\right)\right] d \tau_{1} d \tau_{2} \\
& +\frac{1}{2} \sum_{p_{1}} \sum_{p_{2}} \int_{0}^{1} \int_{0}^{1} \Delta_{p_{2}}\left(S_{p_{1}}+\tau_{1} \Delta_{p_{1}}\right)\left[\tilde{W}\left(\Omega_{p_{1}, p_{2}}\left(\tau_{1}, \tau_{2}\right) v\right) \Delta_{p_{1}}\right] d \tau_{1} d \tau_{2} \tag{2.3.8}
\end{align*}
$$

This is a symmetric operator. We apply (2.2.22) to $\widetilde{W}$. Because of the cut-offs in the argument of $\widetilde{W}$ in (2.3.8), we may write $\Pi_{n} W(v) \Pi_{n^{\prime}}=\Pi_{n} W\left(S\left(n, n^{\prime}\right) v\right) \Pi_{n^{\prime}}$. Consequently, (2.2.22) implies (2.3.7). Note that since $\sigma \geq \sigma_{0}+d+1$, we may take some integer $M>d$, such that $\sigma_{0}+M \leq \sigma$, so that for $v, h_{\ell^{\prime}}$ in $\mathscr{H}^{\sigma}$, the right side of (2.3.7) is bounded from above by $C\left\langle n-n^{\prime}\right\rangle^{-M}$. This shows that $W(v)$ is indeed bounded on $\mathscr{H}^{0}$.

Proof of Proposition 2.3.1. Let $h_{1}$ be in $\mathscr{H}^{+\infty}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{R}^{2}\right)$ and write

$$
\begin{equation*}
D \psi_{2}(v, \omega, \epsilon) \cdot h_{1}=2 \int_{\mathbb{S}^{1} \times \mathbb{T}^{d}}(W(v, \omega, \epsilon) v) h_{1} d t d x+\int_{\mathbb{S}^{1} \times \mathbb{T}^{d}}\left(\left(D W(v, \omega, \epsilon) \cdot h_{1}\right) v\right) v d t d x \tag{2.3.9}
\end{equation*}
$$

Define

$$
\tilde{V}=2 \sum_{n, n^{\prime}} \mathbb{1}_{\left|n-n^{\prime}\right| \leq \frac{1}{10}}\left(|n|+\left|n^{\prime}\right|\right), \Pi_{n} W(v, \omega, \epsilon) \Pi_{n^{\prime}}
$$

In (2.3.7), we can bound $\left\|\partial^{\alpha} S\left(n, n^{\prime}\right) v\right\|_{\mathscr{H} \sigma_{0}}$ by $C\|v\|_{\mathscr{H} \sigma}$ when $|\alpha| \leq M \leq \sigma-\sigma_{0}$, and we can control $\left\|S\left(n, n^{\prime}\right) h_{\ell^{\prime}}\right\|_{\mathscr{H}^{\sigma_{0}}+N_{\ell^{\prime}}}$ by $C\left\|h_{\ell^{\prime}}\right\|_{\mathscr{H}} \sigma_{0}+M$. We obtain that $\tilde{V}$ satisfies (2.1.1), and is thus an element of $\Psi^{0}(0, \sigma, q)$. We show that the remaining terms in (2.3.9) give contributions to the last term in (2.3.3). Set

$$
R_{1}(v, \omega, \epsilon)=2 \sum_{n} \sum_{n^{\prime}} \Pi_{n} W(v, \omega, \epsilon) \Pi_{n^{\prime}} \mathbb{1}_{\left|n-n^{\prime}\right|>\frac{1}{10}\left(|n|+\left|n^{\prime}\right|\right)}
$$

We estimate

$$
\begin{equation*}
\left\|\Pi_{n} \partial_{\omega}^{a_{0}} \partial_{\epsilon}^{a_{1}} \partial_{v}^{\ell} R_{1}(v, \omega, \epsilon) \cdot\left(h_{1}, \ldots, h_{\ell}\right) \Pi_{n^{\prime}}\right\|_{\mathscr{L}(\mathscr{H} 0} \tag{2.3.10}
\end{equation*}
$$

$\operatorname{using}$ (2.3.7) with $M>\sigma-\sigma_{0}$. Since $\left\|S\left(n, n^{\prime}\right) w\right\|_{\mathscr{H} \sigma_{0}+\beta} \leq C\left(1+\inf \left(|n|,\left|n^{\prime}\right|\right)\right)^{\left(\beta+\sigma_{0}-\sigma\right)_{+}}\|w\|_{\mathscr{H} \sigma}$, we get for (2.3.10) the upper bound

$$
C\left(1+|n|+\left|n^{\prime}\right|\right)^{-M}\left(1+\inf \left(|n|,\left|n^{\prime}\right|\right)\right)^{M+\sigma_{0}-\sigma} \prod_{\ell^{\prime}=1}^{\ell}\left\|h_{\ell^{\prime}}\right\|_{\mathscr{H} \sigma} .
$$

Taking $M$ large enough, we deduce the boundedness of $R_{1}(v, \omega, \epsilon)$ and of its derivatives from $\mathscr{H}^{s}$ to $\mathscr{H}^{s+\left(\sigma-\sigma_{0}-d-1\right)}$, for any $s \geq \sigma_{0}$; thus $R_{1} \in \mathscr{R}_{0, \mathbb{R}}^{r}(0, \sigma, q)$.

We treat next the last contribution to (2.3.9), defining an operator $R_{2}(v, \omega, \epsilon)$ by

$$
\begin{equation*}
\int[(D W(v, \omega, \epsilon) \cdot h) v] w d t d x=\int\left[R_{2}(v, \omega, \epsilon) w\right] h d t d x \tag{2.3.11}
\end{equation*}
$$

for any $h, w \in \mathscr{H}^{+\infty}$. On the left side, we decompose the last $v$ as $\sum_{n^{\prime}} \Pi_{n^{\prime}} v$ and $w$ as $\sum_{n} \Pi_{n} w$. We bound the modulus of (2.3.11) by

$$
\begin{equation*}
\sum_{n} \sum_{n^{\prime}}\left\|\Pi_{n} D W(v, \omega, \epsilon) \cdot h \Pi_{n^{\prime}}\right\|_{\mathscr{L}\left(\mathscr{H}^{0}\right)}\left\|\Pi_{n^{\prime}} v\right\|_{\mathscr{H} 0}\left\|\Pi_{n} w\right\|_{\mathscr{H} 0} \tag{2.3.12}
\end{equation*}
$$

To show that $R_{2}(v, \omega, \epsilon)$ is bounded from $\mathscr{H}^{s}$ to $\mathscr{H}^{s+r}$, we bound $\left\|\Pi_{n} w\right\|_{\mathscr{H}^{0}} \leq c_{n}\langle n\rangle^{-s}\|w\|_{\mathscr{H}^{s}}$, for a $\ell^{2}$-sequence $\left(c_{n}\right)_{n}$ and take $h \in \mathcal{H}^{-s-r}$. We use (2.3.7) with $\ell=1$. We have the bound

$$
Q_{N_{0}}^{1}\left(\left(\left\|\partial^{\alpha} S\left(n, n^{\prime}\right) v\right\|_{\mathscr{H}} \sigma_{0}\right)_{\alpha}\right)\left\|S\left(n, n^{\prime}\right) h\right\|_{\mathscr{H} \sigma_{0}+N_{1}} \leq C\left(1+\inf \left(|n|,\left|n^{\prime}\right|\right)\right)^{M+s+r+\sigma_{0}}\|h\|_{\mathscr{H}-s-r}
$$

since $v$ is bounded in $\mathscr{H}^{\sigma}$. Consequently, the general term of (2.3.12) is smaller than

$$
\begin{equation*}
C\left\langle n-n^{\prime}\right\rangle^{-M}\left(1+\inf \left(|n|,\left|n^{\prime}\right|\right)\right)^{M+s+r+\sigma_{0}}\langle n\rangle^{-s} c_{n}\|w\|_{\mathscr{H}}\|h\|_{\mathscr{H}-s-r}\left\langle n^{\prime}\right\rangle^{-\sigma} c_{n^{\prime}}^{\prime}\|v\|_{\mathscr{H} \sigma} \tag{2.3.13}
\end{equation*}
$$

for some $\ell^{2}$-sequence $\left(c_{n^{\prime}}^{\prime}\right)_{n^{\prime}}$. Taking $M=d+1$, and using the value (2.3.2) of $r$ and $s \geq 0, \sigma \geq 0$, one checks that the sum in $n, n^{\prime}$ of (2.3.13) converges. This shows the boundedness of $R_{2}(v, \omega, \epsilon)$ from $\mathscr{H}^{s}$ to $\mathscr{H}^{s+r}$. One treats in the same way $\partial_{\omega}^{a_{0}} \partial_{\epsilon}^{a_{1}} \partial_{v}^{\ell} R_{2}(v, \omega, \epsilon)$. Consequently $R_{2} \in \mathscr{R}_{0, \mathbb{R}}^{r}(0, \sigma, q)$. This concludes the proof of the proposition.

Let us conclude this section writing in complex coordinates the equation we are interested in. By Proposition 2.3.1, Equation (2.3.1) may be written

$$
\begin{equation*}
L_{\omega} v+\epsilon f+\epsilon \tilde{V}(v, \omega, \epsilon) v+\epsilon \tilde{R}(v, \omega, \epsilon) v=0 \tag{2.3.14}
\end{equation*}
$$

We write $v=\left[\begin{array}{l}v_{1} \\ v_{2}\end{array}\right] \in \mathbb{R}^{2}$ and set $u=v_{1}+i v_{2}, U=\left[\begin{array}{l}u \\ \bar{u}\end{array}\right], I^{\prime}=\left[\begin{array}{cc}1 & 0 \\ 0 & -1\end{array}\right]$.
Corollary 2.3.3. Let $q>0, \sigma \geq \sigma_{0}+d+1, r$ given by (2.3.2). There is an element $V(U, \omega, \epsilon)$ in $\Psi^{0}(0, \sigma, q) \otimes \mathcal{M}_{2}(\mathbb{R})$ with $V(U, \omega, \epsilon)^{*}=V(U, \omega, \epsilon)$, there is $R(U, \omega, \epsilon)$ in $\mathscr{R}_{0}^{r}(0, \sigma, q) \otimes M_{2}(\mathbb{R})$ such that (2.3.14) is equivalent to

$$
\begin{equation*}
\left[\left(\omega I^{\prime} D_{t}+(-\Delta+\mu) I\right)+\epsilon V(U, \omega, \epsilon)\right] U=\epsilon R(U, \omega, \epsilon) U+\epsilon f \tag{2.3.15}
\end{equation*}
$$

(where, abusing notation, we write f for $\left[\begin{array}{c}f_{1}+i f_{2} \\ f_{1}-i f_{2}\end{array}\right]$ ).
Proof. Write $\tilde{V}(v, \omega, \epsilon)=\left(\tilde{V}_{i, j}(v, \omega, \epsilon)\right)_{1 \leq i, j \leq 2}, \tilde{R}(v, \omega, \epsilon)=\left(\widetilde{R}_{i, j}(v, \omega, \epsilon)\right)_{1 \leq i, j \leq 2}$ and note that (2.3.14) implies

$$
\begin{align*}
& \left(\omega D_{t}-\Delta+\mu\right) u \\
& \quad=\epsilon\left(f_{1}+i f_{2}\right)-\epsilon V_{11}(U, \omega, \epsilon) u-\epsilon V_{12}(U, \omega, \epsilon) \bar{u}+\epsilon R_{11}(U, \omega, \epsilon) u+\epsilon R_{12}(U, \omega, \epsilon) \bar{u} \tag{2.3.16}
\end{align*}
$$

where we have set

$$
\begin{array}{rlr}
V_{11}=-\frac{1}{2}\left[\tilde{V}_{11}+\tilde{V}_{22}+i\left(\tilde{V}_{21}-\tilde{V}_{12}\right)\right], & V_{12}=-\frac{1}{2}\left[\tilde{V}_{11}-\tilde{V}_{22}+i\left(\tilde{V}_{21}+\tilde{V}_{12}\right)\right], \\
R_{11}=\frac{1}{2}\left[\widetilde{R}_{11}+\widetilde{R}_{22}+i\left(\widetilde{R}_{21}-\widetilde{R}_{12}\right)\right], & R_{12}=\frac{1}{2}\left[\widetilde{R}_{11}-\widetilde{R}_{22}+i\left(\widetilde{R}_{21}+\widetilde{R}_{12}\right)\right] . \tag{2.3.17}
\end{array}
$$

We define $V_{\underline{21}}=\bar{V}_{\tilde{\sim}}$, $V_{22}=\bar{V}_{11}, R_{21}=\bar{R}_{12}, R_{22}=\bar{R}_{11}, V=\left(V_{i j}\right)_{1 \leq i, j \leq 2}, R=\left(R_{i j}\right)_{1 \leq i, j \leq 2}$. Since ${ }^{t} \tilde{V}=\tilde{V}$ and $\tilde{V}=\tilde{V}$, we see that $V^{*}=V$ and (2.3.16), (2.3.17) imply (2.3.15). This concludes the proof.

## 3. Diagonalization of the problem

The goal of this section is to deduce from (2.3.15) a new equation where, up to remainders, $V(U, \omega, \epsilon)$ will be replaced by a block diagonal operator relatively to the decomposition $\mathscr{H}^{0}=\bigoplus_{\alpha} \operatorname{Range}\left(\widetilde{\Pi}_{\alpha}\right)$ coming from (1.2.5). This is the key point that will allow us to avoid using Nash-Moser methods in the construction of the solution performed in Section 4.

### 3.1. Spaces of diagonal and non diagonal operators.

Definition 3.1.1. Let $\sigma \in \mathbb{R}, N \in \mathbb{N}, \sigma \geq \sigma_{0}+d+1+2 N, m \in \mathbb{R}, q>0$.
(i) One denotes by $\Sigma^{m}(N, \sigma, q)$ the space $\Psi^{m}(N, \sigma, q) \otimes \mathcal{M}_{2}(\mathbb{R})$. Abusing notation, we also write $\mathscr{R}_{v}^{r}(N, \sigma, q)$ for $\mathscr{R}_{v}^{r}(N, \sigma, q) \otimes \mathcal{M}_{2}(\mathbb{R})$.
(ii) One denotes by $\Sigma_{\mathrm{D}}^{m}(N, \sigma, q)$ the subspace of $\Sigma^{m}(N, \sigma, q)$ consisting of elements $A(U, \omega, \epsilon)=$ $\left(A_{i j}(U, \omega, \epsilon)\right)_{1 \leq i, j \leq 2}$ such that $A_{12}=A_{21}=0$ and, for any $\alpha, \alpha^{\prime} \in \mathscr{A}$ with $\alpha \neq \alpha^{\prime}$,

$$
\begin{equation*}
\tilde{\Pi}_{\alpha} A_{11}(U, \omega, \epsilon) \tilde{\Pi}_{\alpha^{\prime}} \equiv 0, \quad \tilde{\Pi}_{\alpha} A_{22}(U, \omega, \epsilon) \tilde{\Pi}_{\alpha^{\prime}} \equiv 0 \tag{3.1.1}
\end{equation*}
$$

(iii) One denotes by $\Sigma_{\mathrm{ND}}^{m}(N, \sigma, q)$ the subspace of $\Sigma^{m}(N, \sigma, q)$ made up of elements $A(U, \omega, \epsilon)$ such that, for any $\alpha \in \mathscr{A}$,

$$
\begin{equation*}
\tilde{\Pi}_{\alpha} A_{11}(U, \omega, \epsilon) \tilde{\Pi}_{\alpha} \equiv 0, \quad \tilde{\Pi}_{\alpha} A_{22}(U, \omega, \epsilon) \tilde{\Pi}_{\alpha} \equiv 0 \tag{3.1.2}
\end{equation*}
$$

Clearly, we get a direct sum decomposition $\Sigma^{m}(N, \sigma, q)=\Sigma_{\mathrm{D}}^{m}(N, \sigma, q) \oplus \Sigma_{\mathrm{ND}}^{m}(N, \sigma, q)$.
Definition 3.1.2. Let $\rho \in] 0,1]$.
(i) $\mathscr{L}_{\rho}^{m}(N, \sigma, q)$ denotes the subspace of $\Sigma^{m-\rho}(N, \sigma, q)$ consisting of elements $\left(A_{i j}(U, \omega, \epsilon)\right)_{1 \leq i, j \leq 2}$ that satisfy

$$
\begin{equation*}
A_{11}, A_{22} \in \Psi^{m-\rho}(N, \sigma, q), \quad A_{12}, A_{21} \in \Psi^{m-2}(N, \sigma, q) \tag{3.1.3}
\end{equation*}
$$

(ii) $\mathscr{L}_{\rho}^{\prime m}(N, \sigma, q)$ denotes the subspace of $\Sigma^{m-\rho}(N, \sigma, q)$ consisting of elements $\left(A_{i j}(U, \omega, \epsilon)\right)_{1 \leq i, j \leq 2}$ that satisfy (3.1.3) and

$$
\begin{equation*}
A_{11}^{*}=-A_{11}, \quad A_{22}^{*}=-A_{22}, \quad A_{12}^{*}=A_{21} \tag{3.1.4}
\end{equation*}
$$

Remark. It follows from the definition and from Proposition 2.1.4(ii) that, if $A \in \mathscr{L}_{\rho}^{m_{1}}(N, \sigma, q), B \in$ $\mathscr{L}_{\rho}^{m_{2}}(N, \sigma, q)$ with $\sigma \geq \sigma_{0}+2 N+d+1+\left(m_{1}+m_{2}-2 \rho\right)_{+}$, then $A B$ is the sum of an element of $\mathscr{L}_{\rho}^{m_{1}+m_{2}-\rho}(N, \sigma, q)$ and an element of $\mathscr{R}_{0}^{r}(N, \sigma, q)$ with

$$
r=\sigma-\sigma_{0}-(d+1)-m_{1}-m_{2}+2 \rho-2 N
$$

Proposition 3.1.3. Let $A(U, \omega, \epsilon)$ be a self-adjoint element of $\Sigma_{\mathrm{ND}}^{m}(N, \sigma, q)$. There exist $B(U, \omega, \epsilon)$ in $\mathscr{L}_{\rho}^{\prime m}(N, \sigma, q)$ and $R(U, \omega, \epsilon)$ in $\mathscr{R}_{0}^{r(\sigma, N)-m}(N, \sigma, q)$, with $r(\sigma, N)=\rho\left(\sigma-\sigma_{0}-2 N-d-1\right)$, such that

$$
\begin{equation*}
B(U, \omega, \epsilon)^{*}(\Delta-\mu)+(\Delta-\mu) B(U, \omega, \epsilon)=A(U, \omega, \epsilon)+R(U, \omega, \epsilon) \tag{3.1.5}
\end{equation*}
$$

(where $\rho$ is given by Lemma 1.2.1, for a given $\beta \in] 0, \frac{1}{10}\left[\right.$ ). Moreover $[\Delta, B]$ is in $\Sigma^{m}(N, \sigma, q)$.

Proof. By assumption, we may write

$$
A(U, \omega, \epsilon)=\left[\begin{array}{cc}
a(U, \omega, \epsilon) & b(U, \omega, \epsilon) \\
b(U, \omega, \epsilon)^{*} & c(U, \omega, \epsilon)
\end{array}\right]
$$

with $a^{*}=a, c^{*}=c$, and $\tilde{\Pi}_{\alpha} a \tilde{\Pi}_{\alpha^{\prime}}=0=\tilde{\Pi}_{\alpha} c \tilde{\Pi}_{\alpha^{\prime}}$ if $\alpha, \alpha^{\prime} \in \mathscr{A}$ with $\alpha \neq \alpha^{\prime}$. Write $a=a^{\prime}+a^{\prime \prime}$, with

$$
a^{\prime}=\sum_{n, n^{\prime}} \mathbb{1}_{\left|n-n^{\prime}\right| \leq c\left(|n|+\left|n^{\prime}\right|\right)^{\rho}} \Pi_{n} a \Pi_{n^{\prime}}, \quad a^{\prime \prime}=\sum_{n, n^{\prime}} \mathbb{1}_{\left|n-n^{\prime}\right|>c\left(|n|+\left|n^{\prime}\right|\right)^{\rho}} \Pi_{n} a \Pi_{n^{\prime}}
$$

where $c$ is a small positive constant. Applying (2.1.1) with $M=\sigma-\sigma_{0}-2 N-d-1$, we get

$$
\begin{aligned}
\| \Pi_{n} \partial_{U}^{j} a^{\prime \prime}(U)\left(W_{1}, \ldots, W_{j}\right) & \left.\Pi_{n^{\prime}} \|_{\mathscr{L}(\mathscr{H} 0}\right) \\
& \leq C\left(1+|n|+\left|n^{\prime}\right|\right)^{m-r(\sigma, N)}\left\langle n-n^{\prime}\right\rangle^{-d-1} \mathbb{1}_{\left|n-n^{\prime}\right| \leq \frac{1}{10}\left(|n|+\left|n^{\prime}\right|\right)} \prod_{\ell=1}^{j}\left\|W_{\ell}\right\|_{\mathscr{H} \sigma}
\end{aligned}
$$

which implies a bound of type (2.1.3) for any $s \geq \sigma_{0}$, with $r$ replaced by $r(\sigma, N)-m$. Consequently, $a^{\prime \prime}$ gives a contribution to $R$ in (3.1.5) and, changing notation, we may assume that $a=a^{\prime}$. We do the same for the $c$-contribution, so that we reduce ourselves to $a, c$ verifying that

$$
\begin{equation*}
\Pi_{n} a \Pi_{n^{\prime}}=0 \quad \text { and } \quad \Pi_{n} c \Pi_{n^{\prime}}=0 \quad \text { if }\left|n-n^{\prime}\right|>c\left(|n|+\left|n^{\prime}\right|\right)^{\rho} \tag{3.1.6}
\end{equation*}
$$

We look for

$$
B(U, \omega, \epsilon)=\left[\begin{array}{cc}
a_{1}(U, \omega, \epsilon) & b_{1}(U, \omega, \epsilon) \\
b_{1}(U, \omega, \epsilon)^{*} & c_{1}(U, \omega, \epsilon)
\end{array}\right]
$$

for some $a_{1}, b_{1}, c_{1}$ satisfying $a_{1}^{*}=-a_{1}, c_{1}^{*}=-c_{1}$ such that $A(U, \omega, \epsilon)$ equals the left side of (3.1.5). The latter may be written as

$$
\left[\begin{array}{cc}
{\left[\Delta, a_{1}\right]} & (\Delta-\mu) b_{1}+b_{1}(\Delta-\mu)  \tag{3.1.7}\\
b_{1}^{*}(\Delta-\mu)+(\Delta-\mu) b_{1}^{*} & {\left[\Delta, c_{1}\right]}
\end{array}\right]
$$

Consequently, we have to solve the equations

$$
\begin{equation*}
\left[\Delta, a_{1}\right]=a, \quad(\Delta-\mu) b_{1}+b_{1}(\Delta-\mu)=b, \quad\left[\Delta, c_{1}\right]=c \tag{3.1.8}
\end{equation*}
$$

The first of these is equivalent to

$$
\begin{equation*}
\left(\left|n^{\prime}\right|^{2}-|n|^{2}\right) \Pi_{n} a_{1} \Pi_{n^{\prime}}=\Pi_{n} a \Pi_{n^{\prime}} \text { for any } n, n^{\prime} \in \mathbb{Z}^{d} \tag{3.1.9}
\end{equation*}
$$

Since $A \in \Sigma_{\mathrm{ND}}^{m}(N, \sigma, q)$, Definition 3.1.1(ii) implies that the right side in (3.1.9) vanishes if $n, n^{\prime}$ belong to a same $\Omega_{\alpha}$ of the partition of Lemma 1.2.1. Consequently, we may define

$$
\begin{equation*}
a_{1}(U, \omega, \epsilon)=\sum_{\substack{\alpha, \alpha^{\prime} \in \mathscr{A} \\ \alpha \neq \alpha^{\prime}}} \sum_{n \in \Omega_{\alpha}} \sum_{n^{\prime} \in \Omega_{, \alpha^{\prime}}}\left(\left|n^{\prime}\right|^{2}-|n|^{2}\right)^{-1} \Pi_{n} a(U, \omega, \epsilon) \Pi_{n^{\prime}} \tag{3.1.10}
\end{equation*}
$$

If we use the second lower bound in (1.2.2), Definition 2.1.1, and (3.1.6) with a small enough $c>0$, we see that $a_{1}$ satisfies (2.1.1) with $m$ replaced by $m-\rho$. Thus $a_{1} \in \Psi^{m-\rho}(N, \sigma, q)$, and by (3.1.10) and the fact that $a^{*}=a$, we get $a_{1}^{*}=-a_{1}$. The last equation (3.1.8) is solved in the same way.

We are left with finding $b_{1}(U, \omega, \epsilon)$. The equation giving it is equivalent to

$$
\begin{equation*}
-\left(|n|^{2}+\left|n^{\prime}\right|^{2}+2 \mu\right) \Pi_{n} b_{1} \Pi_{n^{\prime}}=\Pi_{n} b \Pi_{n^{\prime}} \tag{3.1.11}
\end{equation*}
$$

Since $\mu \notin \mathbb{Z}_{-}$by assumption, we may always define $b_{1}$ by division. Coming back to Definition 2.1.1, we see that we get an element of $\Psi^{m-2}(N, \sigma, q)$, which is moreover self-adjoint. This concludes the proof since (3.1.7) shows that by construction $\left[\Delta, a_{1}\right],\left[\Delta, c_{1}\right]$ belong to $\Psi^{m}(N, \sigma, q)$, and since $\Delta b_{1}, b_{1} \Delta$ and their adjoints are in $\Psi^{m}(N, \sigma, q)$.
3.2. Diagonalization theorem. The main result of this subsection is the following one, which gives a reduction for the left side of (2.3.15).

Proposition 3.2.1. Let $r$ be a given positive number and fix an integer $N$ such that $(N+1) \rho \geq r+2$. Let $\sigma \in \mathbb{R}$ satisfy

$$
\begin{equation*}
\sigma \geq \sigma_{0}+2(N+1)+d+1+r / \rho \tag{3.2.1}
\end{equation*}
$$

Let $q>0$ be given. One may find elements $Q_{j}(U, \omega, \epsilon)$ in $\mathscr{L}_{\rho}^{-j \rho}(j, \sigma, q), 0 \leq j \leq N$, elements $V_{\mathrm{D}, j}(U, \omega, \epsilon)$ in $\Sigma_{\mathrm{D}}^{-j \rho}(j, \sigma, q), 0 \leq j \leq N-1$, and an element $R_{1}(U, \omega, \epsilon)$ in $\mathscr{R}_{2}^{r}(N+1, \sigma, q)$, with $C^{1}$ dependence in $(\omega, \epsilon)$, such that if one denotes

$$
Q(U, \omega, \epsilon)=\sum_{j=0}^{N} Q_{j}(U, \omega, \epsilon), V_{\mathrm{D}}(U, \omega, \epsilon)=\sum_{j=0}^{N-1} V_{\mathrm{D}, j}(U, \omega, \epsilon), I^{\prime}=\left[\begin{array}{cc}
1 & 0  \tag{3.2.2}\\
0 & -1
\end{array}\right]
$$

one gets, for any $U \in B_{q}\left(\mathscr{H}^{\sigma}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}^{2}\right)\right)$,

$$
\begin{align*}
&(\operatorname{Id}+\epsilon Q(U, \omega, \epsilon))^{*}\left(\omega I^{\prime} D_{t}+(-\Delta+\mu) I+\epsilon V(U, \omega, \epsilon)\right)(\operatorname{Id}+\epsilon Q(U, \omega, \epsilon)) \\
&=\omega I^{\prime} D_{t}+(-\Delta+\mu) I+\epsilon V_{\mathrm{D}}(U, \omega, \epsilon)-\epsilon R_{1}(U, \omega, \epsilon) \tag{3.2.3}
\end{align*}
$$

We shall prove Proposition 3.2 .1 by constructing recursively $Q_{j}, 0 \leq j \leq N$ so that $Q_{j}$ may be written $Q_{j}=Q_{j}^{\prime}+Q_{j}^{\prime \prime}$ with

$$
\begin{align*}
Q_{j}^{\prime} & \in \mathscr{L}_{\rho}^{\prime-j \rho}(j, \sigma, q), & {\left[\Delta, Q_{j}^{\prime}\right] \in \Sigma^{-j \rho}(j, \sigma, q), } & \\
Q_{j}^{\prime \prime} & \in \mathscr{L}_{\rho}^{-(j+1) \rho}(j, \sigma, q), & {\left[\Delta, Q_{j}^{\prime \prime}\right] \in \Sigma^{-(j+1) \rho}(j, \sigma, q), } & j=0, \ldots, N  \tag{3.2.4}\\
Q_{N}^{\prime \prime} & =0 & &
\end{align*}
$$

We compute first the left side of (3.2.3).
Proposition 3.2.2. Let $r, \sigma, N$ satisfy $(N+1) \rho \geq r+2$ and $\sigma \geq \sigma_{0}+2(N+1)+d+1+r$. Let $Q(U, \omega, \epsilon)=\sum_{j=0}^{N} Q_{j}(U, \omega, \epsilon)$ be given, with $Q_{j}=Q_{j}^{\prime}+Q_{j}^{\prime \prime}$ satisfying (3.2.4).

- There are elements

$$
\begin{equation*}
S_{j}(U, \omega, \epsilon) \in \mathscr{L}_{\rho}^{-(j+1) \rho}(j, \sigma, q), \quad 0 \leq j \leq N-1 \tag{3.2.5}
\end{equation*}
$$

with $\left[\Delta, S_{j}\right] \in \Sigma^{-(j+1) \rho}(j, \sigma, q)$, and $S_{j}$ depending only on $Q_{\ell}^{\prime}(0 \leq \ell \leq j)$ and $Q_{\ell}^{\prime \prime}(0 \leq \ell \leq j-1)$.

- There are elements

$$
V_{j}(U, \omega, \epsilon) \in \Sigma^{-j \rho}(j, \sigma, q), \quad 0 \leq j \leq N
$$

with $\left(V_{j}\right)^{*}=V_{j}, V_{j}$ depending only on $Q_{\ell}(\ell \leq j-1)$.

- There is an element $R \in \mathscr{R}_{2}^{r}(N+1, \sigma, q)$ such that, if we set

$$
\begin{gathered}
V^{N}(U, \omega, \epsilon)=\sum_{j=0}^{N} V_{j}(U, \omega, \epsilon), \quad S^{N}(U, \omega, \epsilon)=\sum_{j=0}^{N-1} S_{j}(U, \omega, \epsilon) \\
Q^{\prime}=\sum_{j=0}^{N} Q_{j}^{\prime}, \quad Q^{\prime \prime}=\sum_{j=0}^{N} Q_{j}^{\prime \prime}, \quad \tilde{L}_{\omega}=\omega I^{\prime} D_{t}+(-\Delta+\mu) I
\end{gathered}
$$

then

$$
\begin{align*}
& (\operatorname{Id}+\epsilon Q)^{*}\left[\widetilde{L}_{\omega}+\epsilon V\right](\operatorname{Id}+\epsilon Q)=\widetilde{L}_{\omega}+\epsilon V^{N}+\epsilon\left[\left(S^{N}\right)^{*} \widetilde{L}_{\omega}+\widetilde{L}_{\omega}\left(S^{N}\right)\right] \\
& +\epsilon\left[Q^{\prime *}(-\Delta+\mu)+(-\Delta+\mu) Q^{\prime}\right]+\epsilon\left[Q^{\prime \prime *} \widetilde{L}_{\omega}+\widetilde{L}_{\omega} Q^{\prime \prime}\right]+\epsilon R \tag{3.2.6}
\end{align*}
$$

Before starting the proof, we compute some commutators.
Lemma 3.2.3. (i) One can find $A_{j} \in \Sigma^{-j \rho}(j-1, \sigma, q)(1 \leq j \leq N)$ depending only on $Q_{\ell}(\ell \leq j-1)$ and satisfying $A_{j}^{*}=A_{j}$, one can find $B_{j} \in \mathscr{L}_{\rho}^{-(j+1) \rho}(j, \sigma, q)(0 \leq j \leq N-1)$ depending only on $Q_{\ell}^{\prime}(\ell \leq j)$ and $Q_{\ell}^{\prime \prime}(\ell \leq j-1)$ and satisfying $\left[\Delta, B_{j}\right] \in \Sigma^{-(j+1) \rho}(j, \sigma, q)$, and one can find $R \in \mathscr{R}_{2}^{r}(N+1, \sigma, q)$, such that, if one sets $A=\sum_{j=1}^{N} A_{j}, B=\sum_{j=0}^{N-1} B_{j}$, then

$$
\begin{equation*}
\left[Q^{*}, \widetilde{L}_{\omega}\right] Q+Q^{*}\left[\widetilde{L}_{\omega}, Q\right]=A+B^{*} \widetilde{L}_{\omega}+\widetilde{L}_{\omega} B+R \tag{3.2.7}
\end{equation*}
$$

(ii) One can find $A_{j}$ as above $(1 \leq j \leq N)$, one can find $B_{j} \in \mathscr{L}_{\rho}^{-(j+1) \rho}(j, \sigma, q)(0 \leq j \leq N-1)$, satisfying $\left[\Delta, B_{j}\right] \in \Sigma^{-(j+1) \rho}(j, \sigma, q)$ and depending only on $Q_{\ell}^{\prime}(\ell \leq j)$ and $Q_{\ell}^{\prime \prime}(\ell \leq j-1)$, and one can find $R \in \mathscr{R}_{2}^{r}(N+1, \sigma, q)$ such that, with the same notation as in $(i)$,

$$
\begin{equation*}
Q^{*} \widetilde{L}_{\omega} Q=A+B^{*} \widetilde{L}_{\omega}+\widetilde{L}_{\omega} B+R \tag{3.2.8}
\end{equation*}
$$

Proof. (i) Write

$$
\begin{aligned}
{\left[\tilde{L}_{\omega}, Q\right] } & =-[\Delta, Q]+\omega\left[I^{\prime} D_{t}, Q\right]=-[\Delta, Q]+\omega I^{\prime}\left[D_{t}, Q\right]+\omega\left[I^{\prime}, Q\right] D_{t} \\
& =-[\Delta, Q]+\omega I^{\prime}\left[D_{t}, Q\right]+\left[I^{\prime}, Q\right] I^{\prime}(\Delta-\mu)+\left[I^{\prime}, Q\right] I^{\prime} \widetilde{L}_{\omega}
\end{aligned}
$$

The left side of (3.2.7) may be written

$$
\begin{array}{r}
-Q^{*}[\Delta, Q]+\omega Q^{*} I^{\prime}\left[D_{t}, Q\right]+Q^{*}\left[I^{\prime}, Q\right] I^{\prime}(\Delta-\mu)-\left[Q^{*}, \Delta\right] Q+\omega\left[Q^{*}, D_{t}\right] I^{\prime} Q+(\Delta-\mu) I^{\prime}\left[Q^{*}, I^{\prime}\right] Q \\
+Q^{*}\left[I^{\prime}, Q\right] I^{\prime} \tilde{L}_{\omega}+\widetilde{L}_{\omega} I^{\prime}\left[Q^{*}, I^{\prime}\right] Q . \tag{3.2.9}
\end{array}
$$

Denote by $\tilde{A}$ the first line of (3.2.9). Then $\tilde{A}$ is self-adjoint and may be written as $\sum_{j=1}^{2 N+2} \tilde{A}_{j}$, where $\tilde{A}_{j}$ is the sum of the following terms:

$$
\begin{array}{ll}
\sum_{\substack{j_{1}+j_{2}=j-1 \\
0 \leq j_{1}, j_{2} \leq N}}\left(-\left[Q_{j_{1}}^{*}, \Delta\right] Q_{j_{2}}-Q_{j_{2}}^{*}\left[\Delta, Q_{j_{1}}\right]\right) & (j \geq 1), \\
\left.\omega \sum_{\substack{j_{1}+j_{2}=j-2 \\
0 \leq j_{1}, j_{2} \leq N}}\left(Q_{j_{1}}^{*} I^{\prime}\left[D_{t}, Q_{j_{2}}\right]+\left[Q_{j_{2}}^{*}, D_{t}\right] I^{\prime} Q_{j_{1}}\right)\right) & (j \geq 2) \\
\sum_{\substack{j_{1}+j_{2}=j-1 \\
0 \leq j_{1}, j_{2} \leq N}}\left(Q_{j_{1}}^{*}\left[I^{\prime}, Q_{j_{2}}\right] I^{\prime}(\Delta-\mu)+(\Delta-\mu) I^{\prime}\left[Q_{j_{2}}^{*}, I^{\prime}\right] Q_{j_{1}}\right) & (j \geq 1) \tag{3.2.11}
\end{array}
$$

Let us check that we may write $\tilde{A}_{j}=A_{j}+R_{1, j}$ with $A_{j}$ in $\Sigma^{-j \rho}(\min (N+1, j-1), \sigma, q)$ and $R_{1, j}$ in $\mathscr{R}_{0}^{r}(\min (N+1, j-1), \sigma, q)$. Since $\mathscr{L}_{\rho}^{-j_{\ell} \rho}\left(j_{\ell}, \sigma, q\right) \subset \Sigma^{-\left(j_{\ell}+1\right) \rho}\left(j_{\ell}, \sigma, q\right)$, it follows from (3.2.4) and from (ii) of Proposition 2.1.4 that the general term in (3.2.10) may be written as a contribution to $A_{j}$ plus a remainder belonging to $\mathscr{R}_{0}^{r_{1}}(\min (N, j-1), \sigma, q)$ with

$$
r_{1}=\sigma-\sigma_{0}-2 N-(d+1)+\left(j_{1}+j_{2}+1\right) \rho \geq r
$$

Moreover these contributions depend only on $Q_{\ell}(\ell \leq j-1)$.
Consider the general term of (3.2.11). The second remark following Definition 2.1.1 implies that $\left[D_{t}, Q_{j_{2}}\right] \in \Sigma^{-\left(j_{2}+1\right) \rho}\left(j_{2}+1, \sigma, q\right)$. Consequently, using again (ii) of Proposition 2.1.4, we may write (3.2.11) as a contribution to $A_{j}$, plus a remainder belonging to $\mathscr{R}_{0}^{r_{1}}(\min (N+1, j-1), \sigma, q)$, depending only on $Q_{\ell}(\ell \leq j-2)$.

Finally, consider (3.2.12). If $C=\left(C_{i j}(U, \omega, \epsilon)\right)_{1 \leq i, j \leq 2}$ is an element of $\mathscr{L}_{\rho}^{m}(N, \sigma, q)$, it follows from (3.1.3) that $\left[I^{\prime}, C\right]=\left[\begin{array}{cc}0 & 2 c_{12} \\ -2 c_{21} & 0\end{array}\right]$ belongs to $\Sigma^{m-2}(N, \sigma, q)$. Hence, the first term in the sum (3.2.12) is given by the composition of an element in $\Sigma^{-\left(j_{1}+1\right) \rho}\left(j_{1}, \sigma, q\right)$ and of an element in $\Sigma^{-j_{2} \rho}\left(j_{2}, \sigma, q\right)$. By applying Proposition 2.1.4 once more, we may write this as a contribution to $A_{j}$ plus a remainder in $\mathscr{R}_{0}^{r}(\min (N, j-1), \sigma, q)$, depending only on $Q_{\ell}(\ell \leq j-1)$. The second term in the argument of the sum (3.2.12) is treated in the same way. This shows that the sum of the first two lines in (3.2.9) contributes to $A+R$ on the right side of (3.2.7), since for $j \geq N+1, A_{j}$ is in $\Sigma^{-(N+1) \rho}(N+1, \sigma, q)$, hence in $\mathscr{R}_{0}^{r}(N+1, \sigma, q)$ by the inequality $(N+1) \rho \geq r$ and the remark after the statement of Definition 2.1.3. Let us show that the last line in (3.2.9) contributes to $B^{*} \widetilde{L}_{\omega}+\widetilde{L}_{\omega} B+R$ in (3.2.7). We have seen above that the fact that $Q_{j}^{\prime} \in \mathscr{L}_{\rho}^{-j \rho}(j, \sigma, q)$ implies $\left[Q_{j}^{\prime}, I^{\prime}\right]=\left[\begin{array}{cc}0 & e_{1} \\ e_{2} & 0\end{array}\right]$ with $e_{\ell} \in \Psi^{-j \rho-2}(j, \sigma, q)$; similarly, $Q_{j}^{\prime \prime} \in \mathscr{L}_{\rho}^{-(j+1) \rho}(j, \sigma, q)$ implies $\left[Q_{j}^{\prime \prime}, I^{\prime}\right]=\left[\begin{array}{cc}0 & e_{1} \\ e_{2} & 0\end{array}\right]$ with $\left.e_{\ell} \in \Psi^{-(j+1) \rho-2}(j, \sigma, q)\right)$. We set
$\tilde{B}_{j}=\sum_{\substack{j_{1}+j_{2}=j \\ 0 \leq j_{1}, j_{2} \leq N}} I^{\prime}\left[Q_{j_{1}}^{\prime *}, I^{\prime}\right] Q_{j_{2}}^{\prime}+\sum_{\substack{j_{1}+j_{2}=j-1 \\ 0 \leq j_{1}, j_{2} \leq N}} I^{\prime}\left(\left[Q_{j_{1}}^{\prime *}, I^{\prime}\right] Q_{j_{2}}^{\prime \prime}+\left[Q_{j_{1}}^{\prime \prime *}, I^{\prime}\right] Q_{j_{2}}^{\prime}\right)+\sum_{\substack{j_{1}+j_{2}=j-2 \\ 0 \leq j_{1}, j_{2} \leq N}} I^{\prime}\left[Q_{j_{1}}^{\prime \prime *}, I^{\prime}\right] Q_{j_{2}}^{\prime \prime}$.
Applying Proposition 2.1.4, we again have a decomposition $\tilde{B}_{j}=B_{j}+R_{j}$, where $B_{j}$ belongs to the class $\mathscr{L}_{\rho}^{-(j+1) \rho}(\min (N, j), \sigma, q)$ (actually, $B_{j}$ is in $\left.\Sigma^{-(j+1) \rho-2}(\min (N, j), \sigma, q)\right)$ and $R_{j}$ belongs to
$\mathscr{R}_{0}^{r+2}(\min (j, N), \sigma, q)$ because of (3.2.1). Moreover, $B_{j}$ depends only on $Q_{\ell}^{\prime}(\ell \leq j)$ and $Q_{\ell}^{\prime \prime}(\ell \leq j-1)$, and by construction, $\left[\Delta, B_{j}\right] \in \Sigma^{-(j+1) \rho}(\min (N, j), \sigma, q)$. For $j \leq N-1$, we get contributions to $B$ and $R$ in (3.2.8), noting that $R_{j} \widetilde{L}_{\omega}, \widetilde{L}_{\omega} R_{j}$ are in $\mathscr{R}_{2}^{r}(N, \sigma, q)$. For $j \geq N, B_{j}$ as well as $R_{j}$ contribute to the remainder in (3.2.7) since $(N+1) \rho \geq r$. This concludes the proof of (i).
(ii) We write

$$
Q^{*} \widetilde{L}_{\omega} Q=\frac{1}{2}\left[Q^{*} Q \widetilde{L}_{\omega}+\widetilde{L}_{\omega} Q^{*} Q\right]+\frac{1}{2}\left[Q^{*}\left[\widetilde{L}_{\omega}, Q\right]+\left[Q^{*}, \widetilde{L}_{\omega}\right] Q\right]
$$

By (i), the last term may be written as a contribution to the right side of (3.2.8). Let us write the first term on the right side under the form $B^{*} \tilde{L}_{\omega}+\widetilde{L}_{\omega} B+R$. We write $Q^{*} Q$ as the sum in $j$ of

$$
\sum_{\substack{j_{1}+j_{2}=j \\ 0 \leq j_{1}, j_{2} \leq N}} Q_{j_{1}}^{\prime *} Q_{j_{2}}^{\prime}+\sum_{\substack{j_{1}+j_{2}=j-1 \\ 0 \leq j_{1}, j_{2} \leq N}}\left(Q_{j_{1}}^{\prime *} Q_{j_{2}}^{\prime \prime}+Q_{j_{1}}^{\prime \prime *} Q_{j_{2}}^{\prime}\right)+\sum_{\substack{j_{1}+j_{2}=j-2 \\ 0 \leq j_{1}, j_{2} \leq N}} Q_{j_{1}}^{\prime \prime *} Q_{j_{2}}^{\prime \prime}
$$

By (3.2.4) and the remark following Definition 3.1.2, this expression may be written as $B_{j}+R_{j}$, where $B_{j} \in \mathscr{L}_{\rho}^{-(j+1) \rho}(\min (N, j), \sigma, q)$ depends only on $Q_{\ell}^{\prime}(\ell \leq j)$ and $Q_{\ell}^{\prime \prime}(\ell \leq j-1),\left[B_{j}, \Delta\right]$ belongs to $\Sigma^{-(j+1) \rho}(\min (N, j), \sigma, q)$, and $R_{j}$ belongs to $\mathscr{R}_{0}^{r_{2}}(\min (N, j), \sigma, q)$, with

$$
r_{2}=\sigma-\sigma_{0}-(d+1)+(j+2) \rho-2 \min (j, N) \geq r+2 .
$$

We obtain contributions to the right side of (3.2.8) when $j \leq N-1$, and to the remainder $R$ when $j \geq N$ since $(N+1) \rho \geq r+2$. This concludes the proof.

Proof of Proposition 3.2.2. We write the left side of (3.2.6) as

$$
\begin{align*}
\widetilde{L}_{\omega}+\epsilon V(U, \omega, \epsilon) & +\epsilon\left[Q^{\prime *}(-\Delta+\mu)+(-\Delta+\mu) Q^{\prime}\right]+\epsilon\left[Q^{\prime \prime *} \widetilde{L}_{\omega}+\widetilde{L}_{\omega} Q^{\prime \prime}\right] \\
& +\epsilon\left[Q^{\prime *} I^{\prime} \omega D_{t}+\omega I^{\prime} D_{t} Q^{\prime}\right]+\epsilon^{2} Q^{*} \widetilde{L}_{\omega} Q+\epsilon^{2}\left[Q^{*} V+V Q\right]+\epsilon^{3} Q^{*} V Q \tag{3.2.13}
\end{align*}
$$

The term $V$ in (3.2.13) contributes to the $V_{0}$ component of $V^{N}$ on the right side of (3.2.6). The first two brackets in (3.2.13) give rise to the last two in (3.2.6). To study the contribution of $Q^{*} \widetilde{L}_{\omega} Q$, we use (3.2.8). The $B_{j}$ component of $B$ on the right side of (3.2.8) contributes to the $S_{j}$ component of $S^{N}$ in (3.2.6). Let us study the third bracket in (3.2.13). By (3.2.4) and Definition 3.1.2, we may write $Q_{j-1}^{\prime}=\left[\begin{array}{ll}a & b \\ b^{*} & c\end{array}\right]$ with $a, c \in \Psi^{-j \rho}(j-1, \sigma, q), b \in \Psi^{-(j-1) \rho-2}(j-1, \sigma, q), a^{*}=-a$, and $c^{*}=-c$. This implies that

$$
Q_{j-1}^{\prime} I^{*} I^{\prime} D_{t}+I^{\prime} D_{t} Q_{j-1}^{\prime}=\left[\begin{array}{cr}
{\left[D_{t}, a\right]} & {\left[D_{t}, b\right]} \\
-\left[D_{t}, b^{*}\right] & -\left[D_{t}, c\right]
\end{array}\right]
$$

is a self-adjoint operator belonging to $\Sigma^{-j \rho}(j, \sigma, q), 1 \leq j \leq N$, by the second remark at the bottom of page 644 . We thus get a contribution to $V_{j}$ in (3.2.6).

Finally, let us check that the last two terms in (3.2.13) may be written as contributions to $V^{N}$ and to $R$ on the right side of (3.2.6). Actually, we may write $Q^{*} V+V Q+\epsilon Q^{*} V Q$ as the sum in $j$ of

$$
\begin{align*}
& Q_{j-1}^{\prime}{ }^{*} V+V Q_{j-1}^{\prime}+Q_{j-2}^{\prime \prime} V+V Q_{j-2}^{\prime \prime} \\
& \quad+\epsilon \sum_{j_{1}+j_{2}=j-2} Q_{j_{1}}^{\prime *} V Q_{j_{2}}^{\prime}+\epsilon \sum_{j_{1}+j_{2}=j-3}\left(Q_{j_{1}}^{\prime \prime *} V Q_{j_{2}}^{\prime}+Q_{j_{1}}^{\prime *} V Q_{j_{2}}^{\prime \prime}\right)+\epsilon \sum_{j_{1}+j_{2}=j-4} Q_{j_{1}}^{\prime \prime *} V Q_{j_{2}}^{\prime \prime} \tag{3.2.14}
\end{align*}
$$

Using that $Q_{j}^{\prime} \in \Sigma^{-(j+1) \rho}(j, \sigma, q), Q_{j}^{\prime \prime} \in \Sigma^{-(j+2) \rho}(j, \sigma, q), V \in \Sigma^{0}(0, \sigma, q)$, we write (3.2.14) as $V_{j}+R_{j}$, where $V_{j}$ depends only on $Q_{\ell}^{\prime}(\ell \leq j-1)$ and $Q_{\ell}^{\prime \prime}(\ell \leq j-2)$ and is in $\Sigma^{-j \rho}(\min (N, j-1), \sigma, q)$ and $R_{j} \in \mathscr{R}_{0}^{r}(N, \sigma, q)$. This concludes the proof.
Proof of Proposition 3.2.1. Let us construct recursively $Q_{j}^{\prime}(0 \leq j \leq N)$ and $Q_{j}^{\prime \prime}(0 \leq j \leq N-1)$ so that the right side of (3.2.6) may be written as the right side of (3.2.3). Assume that $Q_{0}, \ldots, Q_{j-1}$ have been already determined in such a way that the right side of (3.2.6) may be written

$$
\begin{align*}
\widetilde{L}_{\omega}+\epsilon \sum_{j^{\prime}=0}^{j-1} V_{\mathrm{D}, j^{\prime}}+\epsilon \sum_{j^{\prime}=j}^{N-1}\left[S_{j^{\prime}}^{*} \widetilde{L}_{\omega}+\widetilde{L}_{\omega} S_{j^{\prime}}\right] & +\epsilon \sum_{j^{\prime}=j}^{N}\left[Q_{j^{\prime}}^{\prime *}(-\Delta+\mu)+(-\Delta+\mu) Q_{j^{\prime}}^{\prime}\right] \\
& +\epsilon \sum_{j^{\prime}=j}^{N-1}\left[Q_{j^{\prime}}^{\prime \prime} *^{*} \widetilde{L}_{\omega}+\widetilde{L}_{\omega} Q_{j^{\prime}}^{\prime \prime}\right]+\epsilon \sum_{j^{\prime}=j}^{N} V_{j^{\prime}}+\epsilon R \tag{3.2.15}
\end{align*}
$$

Write $V_{j}=\left[\begin{array}{ll}a & b \\ b^{*} & c\end{array}\right]$ with $a, b, c \in \Psi^{-j \rho}(j, \sigma, q), a^{*}=a, c^{*}=c$, and define

$$
V_{\mathrm{D}, j}=\sum_{\alpha \in \mathscr{A}} \widetilde{\Pi}_{\alpha}\left[\begin{array}{ll}
a & 0 \\
0 & c
\end{array}\right] \widetilde{\Pi}_{\alpha}, \quad V_{\mathrm{ND}, j}=V_{j}-V_{\mathrm{D}, j}
$$

Then $V_{\mathrm{D}, j} \in \Sigma_{\mathrm{D}}^{-j \rho}(j, \sigma, q),\left(V_{\mathrm{D}, j}\right)^{*}=V_{\mathrm{D}, j}$ and $V_{\mathrm{ND}, j}$ is in $\Sigma_{\mathrm{ND}}^{-j \rho}(j, \sigma, q),\left(V_{\mathrm{ND}, j}\right)^{*}=V_{\mathrm{ND}, j}$. Moreover $V_{\mathrm{ND}, j}$ depends only on $Q_{\ell}(\ell \leq j-1)$. We apply Proposition 3.1.3 to find $Q_{j}^{\prime} \in \mathscr{L}_{\rho}^{\prime-j \rho}(j, \sigma, q)$ and $R_{j} \in \mathscr{R}_{0}^{r(\sigma, j)+j \rho}(j, \sigma, q)$ such that $Q_{j}^{\prime *}(-\Delta+\mu)+(-\Delta+\mu) Q_{j}^{\prime}=V_{\mathrm{ND}, j}+R_{j}$ and $\left[\Delta, Q_{j}^{\prime}\right]$ is in $\Sigma^{-j \rho}(j, \sigma, q)$. The assumption (3.2.1) on $\sigma$ shows that $R_{j}$ contributes to $R_{1}$ in (3.2.3). Moreover condition (3.2.4) is satisfied by $Q_{j}^{\prime}$, so that we have eliminated the $j$-th component in the fourth and sixth terms of (3.2.15). To eliminate the $j$-th component of the third and fifth terms, we set $Q_{j}^{\prime \prime}=-S_{j}$, $j \leq N-1, Q_{N}^{\prime \prime}=0$. Then condition (3.2.4) is satisfied by $Q_{j}^{\prime \prime}$, and the definition is consistent since $S_{j}$ depends only on $Q_{\ell}^{\prime}(\ell \leq j)$ and $Q_{\ell}^{\prime \prime}(\ell \leq j-1)$. This concludes the proof.

## 4. Iterative scheme

This section will be devoted to the proof of Theorem 1.1.1. We shall construct a solution to (2.3.15) which is equivalent to (1.1.3) - writing this equation under an equivalent form involving the right side of (3.2.3). The first subsection will be devoted to the study of the restriction of the operator $\widetilde{L}_{\omega}+$ $\epsilon V_{\mathrm{D}}(U, \omega, \epsilon)$ to the range of one of the projectors $\tilde{\Pi}_{\alpha}$. We shall show that, for $(\omega, \epsilon)$ outside a subset of small measure, this restriction is invertible. As usual in these problems, the inverse we construct loses derivatives. This will not cause much trouble, since Proposition 3.2.1 allows us to write the equation essentially under the form $\left(\widetilde{L}_{\omega}+\epsilon V_{\mathrm{D}}(U, \omega, \epsilon)\right) W=\epsilon R_{1}(U, \omega, \epsilon) W$ for a new unknown $W$. Since $R_{1}$ is smoothing, it gains enough derivatives to compensate the losses coming from $\left(\widetilde{L}_{\omega}+\epsilon V_{\mathrm{D}}\right)^{-1}$. Because of that, we may construct the solution using a standard iterative scheme.
4.1. Lower bounds for eigenvalues. Let $\left.\left.\gamma_{0} \in\right] 0,1\right], \sigma \in \mathbb{R}, N \in \mathbb{N}, \zeta \in \mathbb{R}_{+}$such that

$$
\sigma \geq \sigma_{0}+\frac{\zeta}{\rho}+2(N+1)+d+1
$$

We denote by $\mathscr{C}^{\sigma}(\zeta)$ the space of functions from $\mathbb{S}^{1} \times \mathbb{T}^{d} \times[1,2] \times\left[0, \gamma_{0}\right]$ to $\mathbb{C}^{2}$,

$$
\begin{equation*}
(t, x, \omega, \epsilon) \longrightarrow U(t, x, \omega, \epsilon) \tag{4.1.1}
\end{equation*}
$$

which are continuous functions of $\omega$ with values in $\mathscr{H}^{\sigma}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}^{2}\right)$ and $C^{1}$ functions of $\omega$ with values in $\mathscr{H}^{\sigma-\zeta-2}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}^{2}\right)$, uniformly in $\epsilon \in\left[0, \gamma_{0}\right]$. We set

$$
\begin{equation*}
\|U\|_{\mathscr{E} \sigma}(\zeta)=\sup _{(\omega, \epsilon) \in[1,2] \times\left[0, \gamma_{0}\right]}\|U(\cdot, \omega, \epsilon)\|_{\mathscr{H} \sigma}+\sup _{(\omega, \epsilon) \in[1,2] \times\left[0, \gamma_{0}\right]}\left\|\partial_{\omega} U(\cdot, \omega, \epsilon)\right\|_{\mathscr{H} \sigma-\zeta-2} \tag{4.1.2}
\end{equation*}
$$

If $\widetilde{\Pi}_{\alpha}$ is the projector of $\mathscr{H}^{0}$ given by (1.2.5), we set $F_{\alpha}=\operatorname{Range}\left(\widetilde{\Pi}_{\alpha}\right), D_{\alpha}=\operatorname{dim} F_{\alpha}$. By (1.2.4) and (1.2.6), $D_{\alpha} \leq C_{1}\langle n(\alpha)\rangle^{\beta d+2}$ for some $C_{1}>0$. We define for $U \in \mathscr{C}^{\sigma}(\zeta), \omega \in[1,2], \epsilon \in\left[0, \gamma_{0}\right]$

$$
\begin{equation*}
A_{\alpha}(\omega ; U, \epsilon)=\tilde{\Pi}_{\alpha}\left(\tilde{L}_{\omega}+\epsilon V_{\mathrm{D}}(U, \omega, \epsilon)\right) \tilde{\Pi}_{\alpha} \tag{4.1.3}
\end{equation*}
$$

This is a self-adjoint operator on $F_{\alpha}$, with $C^{1}$ dependence in $\omega$, since it follows from the expression (3.2.2) of $V_{\mathrm{D}}$, condition (2.1.1) in the definition of $\Psi^{m}(N, \sigma, q)$, the fact that $\partial_{\omega} U \in \mathscr{H}^{\sigma-\zeta-2}$, and the assumption made on $\sigma$, that $\omega \rightarrow \widetilde{\Pi}_{\alpha} V_{\mathrm{D}}(U(t, x, \omega, \epsilon), \omega, \epsilon) \widetilde{\Pi}_{\alpha}$ is $C^{1}$. The main result of this subsection is the following:
Proposition 4.1.1. For any $\mu \in \mathbb{R}-\mathbb{Z}_{-}$and $q>0$, there are $\left.\left.\gamma_{0} \in\right] 0,1\right], C_{0}>0, \mathscr{A}_{0} \subset \mathscr{A}$ a finite subset, and for any $U \in \mathscr{E}^{\sigma}(\zeta)$ with $\|U\|_{\mathscr{C} \sigma}^{(\zeta)} \ll q$, any $\epsilon \in\left[0, \gamma_{0}\right]$, any $\alpha \in \mathscr{A}$, the eigenvalues of $A_{\alpha}$ form a finite family of $C^{1}$ real valued functions of $\omega$, depending on $(U, \epsilon)$,

$$
\begin{equation*}
\omega \rightarrow \lambda_{\ell}^{\alpha}(\omega ; U, \epsilon), 1 \leq \ell \leq D_{\alpha} \tag{4.1.4}
\end{equation*}
$$

satisfying the following properties:
(i) For any $\alpha \in \mathscr{A}$, any $U, U^{\prime} \in \mathscr{H}^{\sigma}$ with $\|U\|_{\mathscr{H} \sigma}<q,\left\|U^{\prime}\right\|_{\mathscr{H} \sigma}<q$, any $\ell \in\left\{1, \ldots, D_{\alpha}\right\}$, any $\in \in\left[0, \gamma_{0}\right]$, and any $\omega \in[1,2]$, there is $\ell^{\prime} \in\left\{1, \ldots, D_{\alpha}\right\}$ such that

$$
\begin{equation*}
\left|\lambda_{\ell}^{\alpha}(\omega ; U, \epsilon)-\lambda_{\ell^{\prime}}^{\alpha}\left(\omega ; U^{\prime}, \epsilon\right)\right| \leq C_{0} \epsilon\left\|U-U^{\prime}\right\|_{\mathscr{H}} \sigma \tag{4.1.5}
\end{equation*}
$$

(ii) For any $a \in \mathscr{A}-\mathscr{A}_{0}$, any $U \in \mathscr{E} \sigma(\zeta)$ with $\|U\|_{\mathscr{E} \sigma}^{(\zeta)}{ }_{\mathscr{C}}<q$, any $\epsilon \in\left[0, \gamma_{0}\right]$, and any $\ell \in\left\{1, \ldots, D_{\alpha}\right\}$, either

$$
\begin{equation*}
C_{0}^{-1}\langle n(\alpha)\rangle^{2} \leq \frac{\partial \lambda_{\ell}^{\alpha}}{\partial \omega}(\omega ; U, \epsilon) \leq C_{0}\langle n(\alpha)\rangle^{2} \quad \text { for any } \omega \text { in }[1,2] \tag{4.1.6}
\end{equation*}
$$

or

$$
\begin{equation*}
-C_{0}\langle n(\alpha)\rangle^{2} \leq \frac{\partial \lambda_{\ell}^{\alpha}}{\partial \omega}(\omega ; U, \epsilon) \leq-C_{0}^{-1}\langle n(\alpha)\rangle^{2} \quad \text { for any } \omega \text { in }[1,2] \tag{4.1.7}
\end{equation*}
$$

(iii) For $\delta \in] 0,1], \epsilon \in\left[0, \gamma_{0}\right], \alpha \in \mathscr{A}$, and $U \in \mathscr{E} \sigma(\zeta)$ with $\|U\|_{\mathscr{E} \sigma(\zeta)}<q$, set

$$
\begin{equation*}
I(\alpha, U, \epsilon, \delta)=\left\{\omega \in[1,2]: \forall \ell \in\left\{1, \ldots, D_{\alpha}\right\},\left|\lambda_{\ell}^{\alpha}(\omega ; U, \epsilon)\right| \geq \delta\langle n(\alpha)\rangle^{-\zeta}\right\} \tag{4.1.8}
\end{equation*}
$$

Then there is a constant $E_{0}$, depending only on the dimension, such that for any $\omega \in I(\alpha, U, \epsilon, \delta)$, $A_{\alpha}(\omega ; U, \epsilon)$ is invertible and

$$
\begin{equation*}
\left\|A_{\alpha}(\omega ; U, \epsilon)^{-1}\right\|_{\mathscr{L}\left(\mathscr{H}^{0}\right)} \leq E_{0} \delta^{-1}\langle n(\alpha)\rangle^{\zeta}, \quad\left\|\partial_{\omega} A_{\alpha}(\omega ; U, \epsilon)^{-1}\right\|_{\mathscr{L}\left(\mathscr{H}^{0}\right)} \leq E_{0} \delta^{-2}\langle n(\alpha)\rangle^{2 \zeta+2} \tag{4.1.9}
\end{equation*}
$$

Proof. The proof of this result is quite classical, and may be found in the references given in the introduction. For completeness, we give it in detail.
(i) By construction, $A_{\alpha}$ is a self-adjoint operator, acting on a space of finite dimension $D_{\alpha}$. Moreover, $A_{\alpha}$ is a $C^{1}$ function of $\omega$ if $U \in \mathscr{E}^{\sigma}(\zeta)$. By a theorem of Rellich (see [Kato 1976, Theorem 6.8], for instance), we know that we may index the eigenvalues of that matrix so that they are $C^{1}$ functions of $\omega$ : $\lambda_{\ell}^{\alpha}(\omega ; U, \epsilon)$, for $1 \leq \ell \leq D_{\alpha}$. Moreover, if $B$ and $B^{\prime}$ are self-adjoint matrices of the same dimension, for any eigenvalue $\lambda_{\ell}(B)$ of $B$ there is an eigenvalue $\lambda_{\ell^{\prime}}\left(B^{\prime}\right)$ of $B^{\prime}$ such that $\left|\lambda_{\ell}(B)-\lambda_{\ell^{\prime}}\left(B^{\prime}\right)\right| \leq\left\|B-B^{\prime}\right\|$. Combining this with the fact that $U \rightarrow A_{\alpha}(\omega ; U, \epsilon)$ is lipschitz with values in $\mathscr{L}\left(\mathscr{H}^{0}\right)$, with lipschitz constant $C \epsilon$, we get (4.1.5).
(ii) Set

$$
\Lambda_{ \pm}^{0}(\alpha)=\left\{ \pm j \omega+|n|^{2}+\mu: j \in \mathbb{N}, n \in \Omega_{\alpha}, K_{0}^{-1}\langle n(\alpha)\rangle^{2} \leq j \leq K_{0}\langle n(\alpha)\rangle^{2}\right\}
$$

so that the spectrum of $\widetilde{\Pi}_{\alpha} \widetilde{L}_{\omega} \widetilde{\Pi}_{\alpha}$ is $\Lambda_{+}^{0}(\alpha) \cup \Lambda_{-}^{0}(\alpha)$. The difference between an eigenvalue in $\Lambda_{+}^{0}(\alpha)$, parametrized by $(j, n)$, and an eigenvalue in $\Lambda_{-}^{0}(\alpha)$, parametrized by $\left(j^{\prime}, n^{\prime}\right)\left(j>0, j^{\prime}<0\right)$ is bounded from below by

$$
\omega\left(j-j^{\prime}\right)+|n|^{2}-\left|n^{\prime}\right|^{2} \geq 2 K_{0}^{-1}\langle n(\alpha)\rangle^{2}-\theta-C\langle n(\alpha)\rangle^{\beta},
$$

by the first estimate (1.2.2), for some $C>0, \beta \in] 0, \frac{1}{10}\left[\right.$. If we take the subset $\mathscr{A}_{0}$ large enough, we get that when $\alpha \in \mathscr{A}-\mathscr{A}_{0}$, the difference between two such eigenvalues is bounded from below by $K_{0}^{-1}\langle n(\alpha)\rangle^{2}$. Consequently, if $0 \leq \epsilon<\gamma_{0}$ small enough, the spectrum of $A_{\alpha}$ may be split in two subsets $\Lambda_{+}(\alpha) \cup \Lambda_{-}(\alpha)$ whose distance is bounded from below by $\frac{1}{2} K_{0}^{-1}\langle n(\alpha)\rangle^{2}$. Let $\Gamma$ be a contour in the complex plane turning once around $\Lambda_{+}^{0}(\alpha)$, of length $O\left(\langle n(\alpha)\rangle^{2}\right)$, such that the distance between $\Gamma$ and the spectrum of $\widetilde{L}_{\omega}^{\alpha}=\widetilde{\Pi}_{\alpha} \widetilde{L}_{\omega} \widetilde{\Pi}_{\alpha}$ is bounded from below by $c\langle n(\alpha)\rangle^{2}$, and such that $\Lambda_{-}^{0}(\alpha)$ is outside $\Gamma$. If $\gamma_{0}$ is small enough, this contour satisfies the same conditions with $\Lambda_{ \pm}^{0}(\alpha)$ replaced by $\Lambda_{ \pm}(\alpha)$ and $\widetilde{L}_{\omega}^{\alpha}$ replaced by $A_{\alpha}$. The spectral projectors $\tilde{\Pi}_{\alpha}^{+}(\omega)$ and $\tilde{\Pi}_{\alpha}^{+, 0}$ associated to the eigenvalues $\Lambda_{+}(\alpha)$ and $\Lambda_{+}^{0}(\alpha)$ of $A_{\alpha}$ and $\widetilde{L}_{\omega}^{\alpha}$, respectively, are given by

$$
\begin{equation*}
\tilde{\Pi}_{\alpha}^{+}(\omega)=\frac{1}{2 i \pi} \int_{\Gamma}\left(\zeta \operatorname{Id}-A_{\alpha}\right)^{-1} d \zeta, \quad \tilde{\Pi}_{\alpha}^{+, 0}=\frac{1}{2 i \pi} \int_{\Gamma}\left(\zeta \operatorname{Id}-\widetilde{L}_{\omega}^{\alpha}\right)^{-1} d \zeta \tag{4.1.10}
\end{equation*}
$$

Note that the second projector is just the orthogonal projector on

$$
\operatorname{Vect}\left\{e^{i(j t+n \cdot x)}: n \in \Omega_{\alpha}, K_{0}^{-1}\langle n(\alpha)\rangle^{2} \leq j \leq K_{0}\langle n(\alpha)\rangle^{2}\right\}
$$

so it is independent of $\omega$. Write

$$
\begin{equation*}
\tilde{\Pi}_{\alpha}^{+}(\omega)-\widetilde{\Pi}_{\alpha}^{+, 0}=\frac{1}{2 i \pi} \int_{\Gamma}\left(\zeta \operatorname{Id}-A_{\alpha}\right)^{-1}\left(A_{\alpha}-\tilde{L}_{\omega}^{\alpha}\right)\left(\zeta \operatorname{Id}-\tilde{L}_{\omega}^{\alpha}\right)^{-1} d \zeta \tag{4.1.11}
\end{equation*}
$$

Using (4.1.3) and the definition of $\widetilde{L}_{\omega}^{\alpha}$ we get

$$
\left\|A_{\alpha}-\widetilde{L}_{\omega}^{\alpha}\right\|_{\mathscr{L}\left(F_{\alpha}\right)}+\left\|\partial_{\omega}\left(A_{\alpha}-\widetilde{L}_{\omega}^{\alpha}\right)\right\|_{\mathscr{L}\left(F_{\alpha}\right)} \leq C \epsilon, \quad\left\|\partial_{\omega} A_{\alpha}\right\|_{\mathscr{L}\left(F_{\alpha}\right)}+\left\|\partial_{\omega} \widetilde{L}_{\omega}^{\alpha}\right\|_{\mathscr{L}\left(F_{\alpha}\right)} \leq C\langle n(\alpha)\rangle^{2} .
$$

Consequently (4.1.11) implies

$$
\begin{aligned}
&\left\|\tilde{\Pi}_{\alpha}^{+}(\omega)-\tilde{\Pi}_{\alpha}^{+, 0}\right\|_{\mathscr{L}\left(F_{\alpha}\right)} \leq C \epsilon\langle n(\alpha)\rangle^{-2} \\
&\left\|\partial_{\omega} \widetilde{\Pi}_{\alpha}^{+}(\omega)\right\|_{\mathscr{L}\left(F_{\alpha}\right)}=\left\|\partial_{\omega}\left(\widetilde{\Pi}_{\alpha}^{+}(\omega)-\widetilde{\Pi}_{\alpha}^{+, 0}\right)\right\|_{\mathscr{L}\left(F_{\alpha}\right)} \leq C \epsilon\langle n(\alpha)\rangle^{-2}
\end{aligned}
$$

Writing

$$
\begin{aligned}
& \widetilde{\Pi}_{\alpha}^{+}(\omega) A_{\alpha} \widetilde{\Pi}_{\alpha}^{+}(\omega)=\left(\widetilde{\Pi}_{\alpha}^{+}(\omega)-\widetilde{\Pi}_{\alpha}^{+, 0}\right) A_{\alpha} \widetilde{\Pi}_{\alpha}^{+}(\omega)+\widetilde{\Pi}_{\alpha}^{+, 0}\left(A_{\alpha}-\widetilde{L}_{\omega}^{\alpha}\right) \widetilde{\Pi}_{\alpha}^{+}(\omega) \\
&+\widetilde{\Pi}_{\alpha}^{+, 0} \widetilde{L}_{\omega}^{\alpha}\left(\widetilde{\Pi}_{\alpha}^{+}(\omega)-\widetilde{\Pi}_{\alpha}^{+, 0}\right)+\widetilde{\Pi}_{\alpha}^{+, 0} \widetilde{L}_{\omega}^{\alpha} \widetilde{\Pi}_{\alpha}^{+, 0}
\end{aligned}
$$

we obtain

$$
\begin{equation*}
\left\|\partial_{\omega}\left[\tilde{\Pi}_{\alpha}^{+}(\omega) A_{\alpha} \tilde{\Pi}_{\alpha}^{+}(\omega)-\tilde{\Pi}_{\alpha}^{+, 0} \tilde{L}_{\omega}^{\alpha} \tilde{\Pi}_{\alpha}^{+, 0}\right]\right\|_{\mathscr{L}\left(F_{\alpha}\right)} \leq C \epsilon \tag{4.1.12}
\end{equation*}
$$

Let $I$ be an interval contained in [1,2] over which one of the eigenvalues $\lambda_{\ell}^{\alpha}(\omega ; U, \epsilon)$ of the matrix $\widetilde{\Pi}_{\alpha}^{+}(\omega) A_{\alpha}(\omega ; U, \epsilon) \widetilde{\Pi}_{\alpha}^{+}(\omega)$ has constant multiplicity $m$, and denote by $P(\omega)$ the associated spectral projector. Then $P(\omega)$ is $C^{1}$ in $\omega \in I$ and satisfies $P(\omega)^{2}=P(\omega)$, whence $P(\omega) P^{\prime}(\omega) P(\omega)=0$. We get therefore for

$$
\lambda_{\ell}^{\alpha}(\omega ; U, \epsilon)=\frac{1}{m} \operatorname{tr}\left[P(\omega) \tilde{\Pi}_{\alpha}^{+}(\omega) A_{\alpha}(\omega ; U, \epsilon) \widetilde{\Pi}_{\alpha}^{+}(\omega) P(\omega)\right]
$$

the equality

$$
\partial_{\omega} \lambda_{\ell}^{\alpha}(\omega ; U, \epsilon)=\frac{1}{m} \operatorname{tr}\left[P(\omega) \partial_{\omega}\left(\tilde{\Pi}_{\alpha}^{+}(\omega) A_{\alpha}(\omega ; U, \epsilon) \tilde{\Pi}_{\alpha}^{+}(\omega)\right) P(\omega)\right]
$$

By (4.1.12), we obtain

$$
\begin{equation*}
\partial_{\omega} \lambda_{\ell}^{\alpha}(\omega ; U, \epsilon)=\frac{1}{m} \operatorname{tr}\left[P(\omega) \partial_{\omega}\left(\tilde{\Pi}_{\alpha}^{+, 0} \tilde{L}_{\omega}^{\alpha} \tilde{\Pi}_{\alpha}^{+, 0}\right) P(\omega)\right]+O(\epsilon) \tag{4.1.13}
\end{equation*}
$$

Since $\widetilde{\Pi}_{\alpha}^{+, 0} \widetilde{L}_{\omega}^{\alpha} \widetilde{\Pi}_{\alpha}^{+, 0}$ is by definition of $\widetilde{L}_{\omega}^{\alpha}$ a diagonal matrix with entries $j \omega+|n|^{2}+\mu, n \in \Omega_{\alpha}$, $K_{0}^{-1}\langle n(\alpha)\rangle^{2} \leq j \leq K_{0}\langle n(\alpha)\rangle^{2}$, we see that (4.1.13) stays between $K_{0}^{-1}\langle n(\alpha)\rangle^{2}-C \epsilon$ and $K_{0}\langle n(\alpha)\rangle^{2}+$ $C \epsilon$. This implies (4.1.6) if $\epsilon \in\left[0, \gamma_{0}\right]$ with $\gamma_{0}$ small enough. The case of eigenvalues corresponding to $\Lambda_{-}(\alpha)$ is treated in a similar way, and gives (4.1.7).
(iii) The first estimate in (4.1.9) follows from the fact that the eigenvalues $\lambda_{\ell}^{\alpha}(\omega ; U, \epsilon)$ of $A_{\alpha}$ satisfy the lower bound given by the definition of (4.1.8). The second estimate is a consequence of the first one and of the fact that $\left\|\partial_{\omega} A_{\alpha}(\omega ; U, \epsilon)\right\|_{\mathscr{L}\left(\mathscr{H}^{0}\right)} \leq C\langle n(\alpha)\rangle^{2}$ by definition of $A_{\alpha}$. This concludes the proof.
4.2. Iterative scheme. This subsection will be devoted to the proof of Theorem 1.1.1, constructing the solution as the limit of an iterative scheme. We fix indices $s, \sigma, N, \zeta, r, \delta$ satisfying the inequalities:

$$
\begin{align*}
& \sigma \geq \sigma_{0}+2(N+1)+d+1+r / \rho, \quad r=\zeta  \tag{4.2.1}\\
& \left.(N+1) \rho \geq r+2, \quad s \geq \sigma+\zeta+2, \quad \delta \in] 0, \delta_{0}\right]
\end{align*}
$$

where $\delta_{0}>0$ will be chosen small enough. We also assume that the parameter $\mu$ is in $\mathbb{R}-\mathbb{Z}_{-}$. We shall solve (2.3.15) when its force term $f$ is given in $\mathscr{H}^{s+\zeta}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}^{2}\right)$. To achieve this goal, the main task will be to construct a sequence $\left(G_{k}, \mathscr{O}_{k}, \psi_{k}, U_{k}, W_{k}\right), k \geq 0$, where $G_{k}, \mathscr{O}_{k}$ will be subsets of
$[1,2] \times\left[0, \delta^{2}\right], \psi_{k}$ will be a real valued function defined on $[1,2] \times\left[0, \delta^{2}\right], U_{k}, W_{k}$ will be functions of $(t, x, \omega, \epsilon) \in \mathbb{S}^{1} \times \mathbb{T}^{d} \times[1,2] \times\left[0, \delta^{2}\right]$ with values in $\mathbb{C}^{2}$. At order $k=0$, we define

$$
\begin{align*}
U_{0} & =W_{0}=0 \\
\mathcal{O}_{0} & =\left\{(\omega, \epsilon) \in[1,2] \times\left[0, \gamma_{0}\right]: \exists \alpha \in \mathscr{A}_{0}, \exists \ell \in\left\{1, \ldots, D_{\alpha}\right\} \text { with }\left|\lambda_{\alpha}^{\ell}(\omega ; 0, \epsilon)\right|<2 \delta\right\}, \tag{4.2.2}
\end{align*}
$$

using the notation of Proposition 4.1.1. For any $\epsilon \in\left[0, \gamma_{0}\right]$, we denote by $\mathbb{O}_{0, \epsilon}$ the $\epsilon$-section of $\mathbb{O}_{0}$ and set

$$
G_{0}=\left\{(\omega, \epsilon) \in[1,2] \times\left[0, \gamma_{0}\right]: d\left(\omega, \mathbb{R}-\mathscr{O}_{0, \epsilon}\right) \geq \frac{\delta}{8 C_{0}^{\prime}}\right\}
$$

where $C_{0}^{\prime}>0$ is a constant such that $\left|\partial_{\omega} \lambda_{\ell}^{\alpha}(\omega ; 0, \epsilon)\right| \leq C_{0}^{\prime}$ for any $\alpha \in \mathscr{A}_{0}$, any $\ell \in\left\{1, \ldots, D_{\alpha}\right\}$, and any $(\omega, \epsilon) \in[1,2] \times\left[0, \gamma_{0}\right]$. Then $\mathcal{O}_{0}$ is an open subset of $[1,2] \times\left[0, \gamma_{0}\right]$ and for any $\epsilon \in\left[0, \gamma_{0}\right], G_{0, \epsilon}$ is a closed subset of $[1,2]$, contained in the open subset $\mathrm{O}_{0, \epsilon}$. By Urysohn's lemma, we may for each fixed $\epsilon$ construct a $C^{1}$ function $\omega \rightarrow \psi_{0}(\omega, \epsilon)$, compactly supported in $\mathbb{O}_{0, \epsilon}$, equal to one on $G_{0, \epsilon}$, such that for any $\omega$ and $\epsilon$ with $0 \leq \psi_{0}(\omega, \epsilon) \leq 1$, we have $\left|\partial_{\omega} \psi_{0}(\omega, \epsilon)\right| \leq C_{1} \delta^{-1}$ for some uniform constant $C_{1}$ depending only on $C_{0}^{\prime}$.

We set

$$
\begin{equation*}
\widetilde{S}_{k}=\sum_{\substack{\alpha \in \mathscr{A} \\\langle n(\alpha)\rangle<2^{k}}} \widetilde{\Pi}_{\alpha}, \quad k \geq 1 \tag{4.2.3}
\end{equation*}
$$

Proposition 4.2.1. There are $\left.\left.\delta_{0} \in\right] 0, \sqrt{\gamma_{0}}\right]$, positive constants $C_{1}, B_{1}, B_{2}$ and, for any $k \geq 1$ and $\delta \in$ ]0, $\delta_{0}$ ], a 5-uple $\left(G_{k}, 0_{k}, \psi_{k}, U_{k}, W_{k}\right)$ satisfying the following conditions:

- $\mathcal{O}_{k}=\left\{(\omega, \epsilon) \in[1,2] \times\left[0, \delta^{2}\right]: \begin{array}{l}\exists \alpha \in \mathscr{A}-\mathscr{A}_{0} \text { and } \ell \in\left\{1, \ldots, D_{\alpha}\right\} \text { such that } \\ 2^{k-1} \leq\langle n(\alpha)\rangle<2^{k},\left|\lambda_{\alpha}^{\ell}\left(\omega ; U_{k-1}, \epsilon\right)\right|<2 \delta 2^{-k \zeta}\end{array}\right\}$,

$$
\begin{equation*}
G_{k}=\left\{(\omega, \epsilon) \in[1,2] \times\left[0, \delta^{2}\right]: d\left(\omega, \mathbb{R}-\widehat{O}_{k, \epsilon}\right) \geq \frac{\delta}{8 C_{0}} 2^{-k(\zeta+2)}\right\} \tag{4.2.4}
\end{equation*}
$$

where $C_{0}$ is the constant in (4.1.6), (4.1.7);

- $\psi_{k}:[1,2] \times\left[0, \delta^{2}\right] \rightarrow[0,1]$ is supported in $\mathbb{O}_{k}$, equal to 1 on $G_{k}$,

$$
\begin{equation*}
C^{1} \text { in } \omega, \text { and satisfies }\left|\partial_{\omega} \psi_{k}(\omega, \epsilon)\right| \leq \frac{C_{1}}{\delta} 2^{k(\zeta+2)} \text { for all }(\omega, \epsilon) \tag{4.2.5}
\end{equation*}
$$

- for any $\epsilon \in\left[0, \delta^{2}\right]$, the function $(t, x, \omega) \rightarrow W_{k}(t, x, \omega, \epsilon)$ is a continuous function of $\omega$ with values in $\mathscr{H}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}^{2}\right)$, which is a $C^{1}$ function of $\omega$ with values in $\mathscr{H}^{s-\zeta-2}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}^{2}\right)$ satisfying

$$
\begin{equation*}
\left\|W_{k}(\cdot, \omega, \epsilon)\right\|_{\mathscr{H} s}+\delta\left\|\partial_{\omega} W_{k}(\cdot, \omega, \epsilon)\right\|_{\mathscr{H} s-\zeta-2} \leq B_{1} \frac{\epsilon}{\delta} \tag{4.2.6}
\end{equation*}
$$

moreover, for any $(\omega, \epsilon) \in[1,2] \times\left[0, \delta^{2}\right]-\bigcup_{k^{\prime}=0}^{k} 0_{k^{\prime}}, W_{k}$ satisfies

$$
\begin{align*}
\left(\widetilde{L}_{\omega}+\epsilon V_{\mathrm{D}}\left(U_{k-1}, \omega, \epsilon\right)\right) W_{k}= & \epsilon \widetilde{S}_{k}\left(\operatorname{Id}+\epsilon Q\left(U_{k-1}, \omega, \epsilon\right)\right)^{*} R\left(U_{k-1}, \omega, \epsilon\right) U_{k-1} \\
& +\epsilon \widetilde{S}_{k}\left[R_{1}\left(U_{k-1}, \omega, \epsilon\right) W_{k-1}\right]+\epsilon \widetilde{S}_{k}\left(\operatorname{Id}+\epsilon Q\left(U_{k-1}, \omega, \epsilon\right)\right)^{*} f \tag{4.2.7}
\end{align*}
$$

where $R$ is defined by (2.3.15) and $Q, V_{\mathrm{D}}, R_{1}$ are defined in (3.2.2), (3.2.3);

- the function $U_{k}$ is defined from $W_{k}$ by

$$
\begin{equation*}
U_{k}(t, x, \omega, \epsilon)=\left(\operatorname{Id}+\epsilon Q\left(U_{k-1}, \omega, \epsilon\right)\right) W_{k} \tag{4.2.8}
\end{equation*}
$$

and it satisfies

$$
\begin{equation*}
\left\|U_{k}-U_{k-1}\right\|_{\mathscr{C}} \leq 2 B_{2} \frac{\epsilon}{\delta} 2^{-k \zeta}, \quad\left\|U_{k}(\cdot, \omega, \epsilon)\right\|_{\mathscr{H} s}+\delta\left\|\partial_{\omega} U_{k}(\cdot, \omega, \epsilon)\right\|_{\mathscr{H} s-\zeta-2} \leq B_{2} \frac{\epsilon}{\delta} \tag{4.2.9}
\end{equation*}
$$

moreover,

$$
\begin{equation*}
\left\|W_{k}-W_{k-1}\right\|_{\mathscr{H} \sigma} \leq B_{2} \frac{\epsilon}{\delta} 2^{-k \zeta} \tag{4.2.10}
\end{equation*}
$$

Remark. Since we assume $\epsilon \leq \delta^{2}$, the second inequality in (4.2.9) implies, with the notation introduced in (4.1.2), the uniform bound

$$
\begin{equation*}
\left\|U_{k}\right\|_{\mathscr{C} \sigma}(\zeta)<q \tag{4.2.11}
\end{equation*}
$$

for some $q$.
Let us write the equation for $U_{k}$ following from (4.2.8) and (4.2.7). Because of the uniform estimate (4.2.11) for $U_{k-1}$, if $0 \leq \delta \leq \delta_{0}$ with $\delta_{0}$ small enough, $\left(\operatorname{Id}+\epsilon Q\left(U_{k-1}, \omega, \epsilon\right)\right)^{*}$ is invertible for any $(\omega, \epsilon) \in[1,2] \times\left[0, \delta^{2}\right]$. If we write

$$
\left(\tilde{L}_{\omega}+\epsilon V\left(U_{k-1}, \omega, \epsilon\right)\right) U_{k}=\left(\tilde{L}_{\omega}+\epsilon V\left(U_{k-1}, \omega, \epsilon\right)\right)\left(\operatorname{Id}+\epsilon Q\left(U_{k-1}, \omega, \epsilon\right)\right) W_{k}
$$

and use (3.2.3) multiplied on the left by $\left(\operatorname{Id}+\epsilon Q\left(U_{k-1}, \omega, \epsilon\right)^{*}\right)^{-1}$ and (4.2.7), we get

$$
\begin{align*}
& \left(\widetilde{L}_{\omega}+\epsilon V\left(U_{k-1}, \omega, \epsilon\right)\right) U_{k}=\epsilon\left(\operatorname{Id}+\epsilon Q\left(U_{k-1}, \omega, \epsilon\right)^{*}\right)^{-1}\left[\widetilde{S}_{k}\left(\mathrm{Id}+\epsilon Q\left(U_{k-1}, \omega, \epsilon\right)^{*}\right) R\left(U_{k-1}, \omega, \epsilon\right) U_{k-1}\right. \\
& \left.\quad+\widetilde{S}_{k} R_{1}\left(U_{k-1}, \omega, \epsilon\right) W_{k-1}+\widetilde{S}_{k}\left(\operatorname{Id}+\epsilon Q\left(U_{k-1}, \omega, \epsilon\right)^{*}\right) f-R_{1}\left(U_{k-1}, \omega, \epsilon\right) W_{k}\right] \tag{4.2.12}
\end{align*}
$$

for any $(\omega, \epsilon) \in[1,2] \times\left[0, \delta^{2}\right]-\bigcup_{k^{\prime}=0}^{k} O_{k^{\prime}}, \delta \in\left[0, \delta_{0}\right]$.
Proof of Proposition 4.2.1. We assume that $\left(G_{k}, 0_{k}, \psi_{k}, U_{k}, W_{k}\right)$ have been constructed satisfying (4.2.4) to (4.2.9), and shall construct these data at rank $k+1$, if $\delta_{0}$ is small enough and the constants $C_{1}, B_{1}, B_{2}$ are large enough.

The sets $\mathbb{O}_{k+1}, G_{k+1}$ are defined by (4.2.4) at rank $k+1$ as soon as $U_{k}$ is given. Then for fixed $\epsilon, G_{k+1, \epsilon}$ is a compact subset of the open set $0_{k+1, \epsilon}$, whose distance to the complement of $\mathbb{O}_{k+1, \epsilon}$ is bounded from below by $\frac{\delta}{8 C_{0}} 2^{-(k+1)(\zeta+2)}$. We may construct by Urysohn's lemma a function $\psi_{k+1}$ satisfying (4.2.5) at rank $k+1$. Let us construct $W_{k+1}$ for $(\omega, \epsilon) \in[1,2] \times\left[0, \delta^{2}\right]-\bigcup_{k^{\prime}=0}^{k+1} G_{k^{\prime}}$. Since $V_{\mathrm{D}}\left(U_{k}, \omega, \epsilon\right)$ is by construction a block-diagonal operator, we may write (4.2.7) at rank $k+1$ as the system of equations

$$
\begin{align*}
&\left(\tilde{L}_{\omega}+\epsilon V_{\mathrm{D}}\left(U_{k}, \omega, \epsilon\right)\right) \widetilde{\Pi}_{\alpha} W_{k+1}=\epsilon \widetilde{\Pi}_{\alpha} \widetilde{S}_{k+1}\left(\operatorname{Id}+\epsilon Q\left(U_{k}, \omega, \epsilon\right)^{*}\right) R\left(U_{k}, \omega, \epsilon\right) U_{k} \\
&+\epsilon \widetilde{\Pi}_{\alpha} \widetilde{S}_{k+1} R_{1}\left(U_{k}, \omega, \epsilon\right) W_{k}+\epsilon \widetilde{\Pi}_{\alpha} \widetilde{S}_{k+1}\left(\operatorname{Id}+\epsilon Q\left(U_{k}, \omega, \epsilon\right)^{*}\right) f \tag{4.2.13}
\end{align*}
$$

for any $\alpha \in \mathscr{A}$. If $\langle n(\alpha)\rangle \geq 2^{k+1}$, the right side of (4.2.13) vanishes by definition of $\tilde{S}_{k+1}$, so that we may set in this case $\widetilde{\Pi}_{\alpha} W_{k+1}=0$ by definition. Let us solve (4.2.13) for those $\alpha$ satisfying $\langle n(\alpha)\rangle<2^{k+1}$. We shall apply Proposition 4.1.1, using the following lemma:

Lemma 4.2.2. There is $\left.\left.\delta_{0} \in\right] 0,1\right]$, depending only on the constants $B_{1}, B_{2}$, such that for any $k \geq 0$, any $k^{\prime} \in\{1, \ldots, k+1\}$, any $\delta \in\left[0, \delta_{0}\right]$, any $\in \in\left[0, \delta^{2}\right]$, and any $\alpha \in \mathscr{A}-\mathscr{A}_{0}$ with $2^{k^{\prime}} \leq\langle n(\alpha)\rangle<2^{k^{\prime}+1}$,

$$
\begin{equation*}
[1,2]-G_{k^{\prime}, \epsilon} \subset I\left(\alpha, U_{k}, \epsilon, \delta\right) \tag{4.2.14}
\end{equation*}
$$

where $I(\cdot)$ is defined by (4.1.8). The same conclusion holds when $k^{\prime}=0, \alpha \in \mathscr{A}_{0}$.
Proof. Consider first the case $k^{\prime} \neq 0$. Let $\omega \in[1,2]-\mathcal{O}_{k^{\prime}, \epsilon}$. Take $\ell \in\left\{1, \ldots, D_{\alpha}\right\}$. By (i) of Proposition 4.1.1 applied to $\left(U, U^{\prime}\right)=\left(U_{k}, U_{k^{\prime}-1}\right)$, there is $\ell^{\prime} \in\left\{1, \ldots, D_{\alpha}\right\}$ such that

$$
\begin{equation*}
\left|\lambda_{\ell}^{\alpha}\left(\omega ; U_{k}, \epsilon\right)\right| \geq\left|\lambda_{\ell^{\prime}}^{\alpha}\left(\omega ; U_{k^{\prime}-1}, \epsilon\right)\right|-C_{0} \epsilon\left\|U_{k}-U_{k^{\prime}-1}\right\|_{\mathscr{H} \sigma} \geq 2 \delta 2^{-k^{\prime} \zeta}-2 C_{0} B_{2} \frac{\epsilon^{2}}{\delta} \frac{2^{-k^{\prime} \zeta}}{1-2^{-\zeta}} \tag{4.2.15}
\end{equation*}
$$

where the second lower bound follows from the definition (4.2.4) of $\mathcal{O}_{k^{\prime}}$ and from (4.2.9). Since $\epsilon \leq \delta^{2}$, we obtain the lower bound

$$
\begin{equation*}
\left|\lambda_{\ell}^{\alpha}\left(\omega ; U_{k}, \epsilon\right)\right| \geq \frac{3}{2} \delta 2^{-k^{\prime} \zeta} \tag{4.2.16}
\end{equation*}
$$

if $\omega \in[1,2]-\mathscr{O}_{k^{\prime}, \epsilon}$ and $\delta \in\left[0, \delta_{0}\right]$ with $\delta_{0}$ small enough. If $\omega \in \mathbb{O}_{k^{\prime}, \epsilon}-G_{k^{\prime}, \epsilon}$, we take $\tilde{\omega} \in[1,2]-\mathcal{O}_{k^{\prime}, \epsilon}$ with $|\omega-\tilde{\omega}|<\frac{\delta}{8 C_{0}} 2^{-k^{\prime}(\zeta+2)}$. By (4.1.6), (4.1.7), we know that for any $U \in \mathscr{E}^{\sigma}(\zeta)$ with $\|U\|_{\mathscr{E} \sigma}(\zeta)<q$, any $\alpha \in \mathscr{A}-\mathscr{A}_{0}$, any $\ell \in\left\{1, \ldots, D_{\alpha}\right\}$,

$$
\sup _{\omega^{\prime} \in[1,2]}\left|\partial_{\omega} \lambda_{\ell}^{\alpha}\left(\omega^{\prime} ; U, \epsilon\right)\right| \leq C_{0}\langle n(\alpha)\rangle^{2}
$$

Enlarging $C_{0}$, we may assume that this inequality is also valid when $\alpha \in \mathscr{A}_{0}$. By condition (4.2.11), we may apply it when $U=U_{k}$. Using (4.2.16), we get since $2^{2 k^{\prime}} \leq\langle n(\alpha)\rangle^{2}<2^{2\left(k^{\prime}+1\right)}$

$$
\left|\lambda_{\ell}^{\alpha}\left(\omega ; U_{k}, \epsilon\right)\right| \geq\left|\lambda_{\ell}^{\alpha}\left(\tilde{\omega} ; U_{k}, \epsilon\right)\right|-C_{0}\langle n(\alpha)\rangle^{2}|\omega-\tilde{\omega}| \geq \delta 2^{-k^{\prime} \zeta} \geq \delta\langle n(\alpha)\rangle^{-\zeta}
$$

When $k^{\prime}=0$, we argue in the same way, taking in (4.2.15) $U_{k^{\prime}-1}=0$. This shows that $\omega$ belongs to $I\left(\alpha, U_{k}, \epsilon, \delta\right)$.

To solve (4.2.13), we shall need, in addition to the preceding lemma, estimates for its right side. Set $H_{k+1}\left(U_{k}, W_{k}\right)=\widetilde{S}_{k+1}\left(\operatorname{Id}+\epsilon Q\left(U_{k}, \omega, \epsilon\right)^{*}\right) R\left(U_{k}, \omega, \epsilon\right) U_{k}$

$$
\begin{equation*}
+\widetilde{S}_{k+1} R_{1}\left(U_{k}, \omega, \epsilon\right) W_{k}+\widetilde{S}_{k+1}\left(\operatorname{Id}+\epsilon Q\left(U_{k}, \omega, \epsilon\right)^{*}\right) f \tag{4.2.17}
\end{equation*}
$$

Lemma 4.2.3. There is a constant $C>0$, depending on $q$ in (4.2.11) but independent of $k$, such that for any $\omega \in[1,2]$, any $\epsilon \in\left[0, \delta^{2}\right]$, and any $\left.\left.\delta \in\right] 0, \delta_{0}\right]$,

$$
\begin{align*}
&\left\|H_{k+1}\left(U_{k}, W_{k}\right)\right\|_{\mathscr{H} s+\zeta} \leq C[\|\left.U_{k}(\cdot, \omega, \epsilon)\left\|_{\mathscr{H}^{s}}+\right\| W_{k}(\cdot, \omega, \epsilon) \|_{\mathscr{H} s}\right]+(1+C \epsilon)\|f\|_{\mathscr{H} s+\zeta}  \tag{4.2.18}\\
&\left\|\partial_{\omega} H_{k+1}\left(U_{k}, W_{k}\right)\right\|_{\mathscr{H}} s-2 \leq C[\| U_{k}(\cdot, \omega, \epsilon)\left\|_{\mathscr{H}^{s}}+\right\| \partial_{\omega} U_{k}(\cdot, \omega, \epsilon) \|_{\mathscr{H}} s-\zeta-2 \\
&\left.+\left\|W_{k}(\cdot, \omega, \epsilon)\right\|_{\mathscr{H} s}+\left\|\partial_{\omega} W_{k}(\cdot, \omega, \epsilon)\right\|_{\mathscr{H} s-\zeta-2}+\epsilon\|f\|_{\mathscr{H} s-2}\right]  \tag{4.2.19}\\
&\left\|H_{k+1}\left(U_{k}, W_{k}\right)-H_{k}\left(U_{k-1}, W_{k-1}\right)\right\|_{\mathscr{H} \sigma+\zeta} \leq C\left[\left\|U_{k}-U_{k-1}\right\|_{\mathscr{H} \sigma}+\left\|W_{k}-W_{k-1}\right\|_{\mathscr{H} \sigma}\right] \\
&+2^{-k \zeta}\left[C\left(\left\|U_{k}\right\|_{\mathscr{H} \sigma+\zeta}+\left\|W_{k}\right\|_{\mathscr{H} \sigma+\zeta}\right)+(1+C \epsilon)\|f\|_{\mathscr{H} \sigma+2 \zeta}\right] . \tag{4.2.20}
\end{align*}
$$

Proof. The operators $R$ and $R_{1}$ belong to $\mathscr{R}_{2}^{r}(N+1, \sigma, q)$ with $r=\zeta$. By Definition 2.1.3, and because of the assumption (4.2.1) on the indices, they are bounded from $\mathscr{H}^{s}$ to $\mathscr{H}^{s+\zeta}$. Moreover, $Q\left(U_{k}, \omega, \epsilon\right)^{*}$ is in $\Psi^{0}(N, \sigma, q) \otimes \mathcal{M}_{2}(\mathbb{R})$, so is bounded on any $\mathscr{H}^{s}$-space by Lemma 2.1.2. This gives (4.2.18).

To obtain (4.2.19), one has to study the boundedness properties of

$$
\begin{align*}
\partial_{\omega}\left[Q\left(U_{k}, \omega, \epsilon\right)\right] & =\partial_{U} Q(\cdot, \omega, \epsilon) \cdot\left(\partial_{\omega} U_{k}\right)+\partial_{\omega} Q\left(U_{k}, \omega, \epsilon\right),  \tag{4.2.21a}\\
\partial_{\omega}\left[R\left(U_{k}, \omega, \epsilon\right)\right] & =\partial_{U} R(\cdot, \omega, \epsilon) \cdot\left(\partial_{\omega} U_{k}\right)+\partial_{\omega} R\left(U_{k}, \omega, \epsilon\right),  \tag{4.2.21b}\\
\partial_{\omega}\left[R_{1}\left(U_{k}, \omega, \epsilon\right)\right] & =\partial_{U} R_{1}(\cdot, \omega, \epsilon) \cdot\left(\partial_{\omega} U_{k}\right)+\partial_{\omega} R_{1}\left(U_{k}, \omega, \epsilon\right) \tag{4.2.21c}
\end{align*}
$$

By (2.1.2), the inequalities in (4.2.1), and the fact that, by (4.2.11), $\partial_{\omega} U_{k}$ is uniformly bounded in $\mathscr{H}^{s-\zeta-2} \subset \mathscr{H}^{\sigma}$, we see that the operator in (4.2.21a) is bounded on any space $\mathscr{H}^{s^{\prime}}$. By (2.1.3), and the assumption $s \geq \sigma+\zeta+2$ in (4.2.1), we see in the same way that (4.2.21b) and (4.2.21c) give bounded operators from $\mathscr{H}^{s-\zeta-2}$ to $\mathscr{H}^{s-2}$ and from $\mathscr{H}^{s}$ to $\mathscr{H}^{s+\zeta . ~ T h i s ~ g i v e s ~ e s t i m a t e ~(4.2 .19) . ~}$

To prove (4.2.20), let us write the difference $H_{k+1}\left(U_{k}, W_{k}\right)-H_{k}\left(U_{k-1}, W_{k-1}\right)$ from the quantities

$$
\left.\begin{array}{r}
\left(\widetilde{S}_{k+1}-\widetilde{S}_{k}\right)\left(\operatorname{Id}+\epsilon Q\left(U_{k}, \omega, \epsilon\right)^{*}\right) R\left(U_{k}, \omega, \epsilon\right) U_{k}, \\
\left(\widetilde{S}_{k+1}-\widetilde{S}_{k}\right) R_{1}\left(U_{k}, \omega, \epsilon\right) W_{k}, \\
\left(\widetilde{S}_{k+1}-\widetilde{S}_{k}\right)\left(\operatorname{Id}+\epsilon Q\left(U_{k}, \omega, \epsilon\right)^{*}\right) f,
\end{array}\right\}
$$

By (4.2.6) and (4.2.9), $U_{k}$ and $W_{k}$ stay in a bounded subset of $\mathscr{H}^{\sigma}$ and $R, R_{1}$ act from $H^{\sigma+\zeta}$ to $H^{\sigma+2 \zeta}$. Using the cut-off $\widetilde{S}_{k+1}-\widetilde{S}_{k}$, we see that the $\mathscr{H}{ }^{\sigma+\zeta}$ norm of (4.2.22) is bounded from above by the last term in the right side of (4.2.20).

By (2.1.3), the $\mathscr{L}\left(\mathscr{H}^{\sigma}, \mathscr{H}^{\sigma+\zeta}\right)$ operator norm of $R\left(U_{k}, \omega, \epsilon\right)-R\left(U_{k-1}, \omega, \epsilon\right)$ and of $R_{1}\left(U_{k}, \omega, \epsilon\right)-$ $R_{1}\left(U_{k-1}, \omega, \epsilon\right)$ is bounded from above by $C\left\|U_{k}-U_{k-1}\right\|_{\mathscr{H}} \sigma$. By (2.1.2), the $\mathscr{L}\left(\mathscr{H}^{\sigma+\zeta}\right.$, $\left.\mathscr{H}^{\sigma+\zeta}\right)$-norm of $Q\left(U_{k}, \omega, \epsilon\right)^{*}-Q\left(U_{k-1}, \omega, \epsilon\right)^{*}$ is bounded by the same quantity. This shows that the $\mathscr{H}{ }^{\sigma+\zeta}$ norm of (4.2.23) is bounded from above by the right side of (4.2.20).

Finally, (4.2.24) is trivially estimated. This concludes the proof.
We continue with the proof of Proposition 4.2.1. We have seen that $\widetilde{\Pi}_{\alpha} W_{k+1}$ is a solution to (4.2.13). Let $k^{\prime} \in\{1, \ldots, k+1\}$ and $\alpha \in \mathscr{A}-\mathscr{A}_{0}$ such that $2^{k^{\prime}} \leq\langle n(\alpha)\rangle<2^{k^{\prime}+1}$, or $k^{\prime}=0, \alpha \in \mathscr{A}_{0}$. Let $\omega \in[1,2]-G_{k^{\prime}, \epsilon}$. By Lemma 4.2.2 and Proposition 4.1.1, the operator $A_{\alpha}\left(\omega ; U_{k}, \epsilon\right)$ is invertible, ant its inverse satisfies estimates (4.1.9). For such $\omega$, we may write (4.2.13) as

$$
\begin{equation*}
\tilde{\Pi}_{\alpha} W_{k+1}=\epsilon A_{\alpha}\left(\omega ; U_{k}, \epsilon\right)^{-1} \tilde{\Pi}_{\alpha} H_{k+1}\left(U_{k}, W_{k}\right) \tag{4.2.25}
\end{equation*}
$$

Applying (4.1.9) we obtain, for any $k^{\prime} \in\{1, \ldots, k+1\}$, any $\alpha \in \mathscr{A}-\mathscr{A}_{0}$ with $2^{k^{\prime}} \leq\langle n(\alpha)\rangle<2^{k^{\prime}+1}$, and any $(\omega, \epsilon) \in[1,2] \times\left[0, \delta^{2}\right]-G_{k^{\prime}}$ (and also any $\alpha \in \mathscr{A}_{0}$ and $(\omega, \epsilon) \in[1,2] \times\left[0, \delta^{2}\right]-G_{0}$ ), the bound

$$
\begin{equation*}
\left\|\tilde{\Pi}_{\alpha} W_{k+1}(\cdot, \omega, \epsilon)\right\|_{\mathscr{H} s} \leq E_{0} \frac{\epsilon}{\delta}\left\|\tilde{\Pi}_{\alpha} H_{k+1}\left(U_{k}, W_{k}\right)(\cdot, \omega, \epsilon)\right\|_{\mathscr{H} s+\zeta} \tag{4.2.26}
\end{equation*}
$$

In the same way, one gets the estimate
$\left\|\tilde{\Pi}_{\alpha} \partial_{\omega} W_{k+1}(\cdot, \omega, \epsilon)\right\|_{\mathscr{H} s-\zeta-2}$

$$
\begin{equation*}
\leq E_{0} \frac{\epsilon}{\delta}\left\|\tilde{\Pi}_{\alpha} \partial_{\omega} H_{k+1}\left(U_{k}, W_{k}\right)(\cdot, \omega, \epsilon)\right\|_{\mathscr{H} s-2}+E_{0} \frac{\epsilon}{\delta^{2}}\left\|\tilde{\Pi}_{\alpha} H_{k+1}\left(U_{k}, W_{k}\right)(\cdot, \omega, \epsilon)\right\|_{\mathscr{H} s+\zeta} \tag{4.2.27}
\end{equation*}
$$

We define $W_{k+1}(t, x, \omega, \epsilon)$ for any value of $(\omega, \epsilon)$ in $[1,2] \times\left[0, \delta^{2}\right]$ from (4.2.25) by setting

$$
\begin{align*}
& W_{k+1}(t, x, \omega, \epsilon)=\sum_{k^{\prime}=1}^{k+1} \sum_{\substack{\alpha \in \mathscr{A}-\mathscr{A}_{0} \\
2^{k^{\prime} \leq\langle n(\alpha)\rangle<2^{k^{\prime}+1}}}}\left(1-\psi_{k^{\prime}}(\omega, \epsilon)\right) \widetilde{\Pi}_{\alpha} W_{k+1}(t, x, \omega, \epsilon)  \tag{4.2.28}\\
&+\sum_{\alpha \in \mathscr{A}_{0}}\left(1-\psi_{0}\right)(\omega, \epsilon) \widetilde{\Pi}_{\alpha} W_{k+1}(t, x, \omega, \epsilon)
\end{align*}
$$

Note that the right side is well defined since (4.2.25) determines $\widetilde{\Pi}_{\alpha} W_{k+1}(\cdot, \omega, \epsilon)$ on the support of $1-\psi_{k^{\prime}}$ when $\left(\alpha, k^{\prime}\right)$ satisfy the conditions in the summation.

We combine (4.2.28), (4.2.26) and (4.2.18). Taking into account (4.2.6) and (4.2.9), we get

$$
\begin{equation*}
\left\|W_{k+1}(\cdot, \omega, \epsilon)\right\|_{\mathscr{H}} \leq E_{0} \frac{\epsilon}{\delta}\left(C\left(B_{1}+B_{2}\right) \frac{\epsilon}{\delta}+\|f\|_{\mathscr{H} s+\zeta}(1+C \epsilon)\right) \tag{4.2.29}
\end{equation*}
$$

To bound the $\partial_{\omega}$-derivative, we use that by (4.2.5)

$$
\left\|\partial_{\omega} \psi_{k^{\prime}} \tilde{\Pi}_{\alpha} W_{k+1}\right\|_{\mathscr{H} s-\zeta-2} \leq \frac{C_{1}}{\delta}\left\|\tilde{\Pi}_{\alpha} W_{k+1}\right\|_{\mathscr{H} s}
$$

when $2^{k^{\prime}} \leq\langle n(\alpha)\rangle<2^{k^{\prime}+1}, \alpha \in \mathscr{A}-\mathscr{A}_{0}$ if $k^{\prime} \neq 0$, and when $\alpha \in \mathscr{A}_{0}$ if $k^{\prime}=0$. We apply this inequality together with (4.2.28), (4.2.27), (4.2.18), (4.2.19) and the uniform bounds (4.2.6), (4.2.9), to get $\left\|\partial_{\omega} W_{k+1}(\cdot, \omega, \epsilon)\right\|_{\mathscr{H} s-\zeta-2} \leq E_{0} \frac{\epsilon}{\delta}\left(C\left(B_{1}+B_{2}\right) \frac{\epsilon}{\delta^{2}}+C \epsilon\|f\|_{\mathscr{H} s-2}\right)$
$+E_{0} \frac{\epsilon}{\delta^{2}}\left(C\left(B_{1}+B_{2}\right) \frac{\epsilon}{\delta}+(1+C \epsilon)\|f\|_{\mathscr{H} s+\zeta}\right)+E_{0} C_{1} \frac{\epsilon}{\delta^{2}}\left(C\left(B_{1}+B_{2}\right) \frac{\epsilon}{\delta}+(1+C \epsilon)\|f\|_{\mathscr{H} s+\zeta}\right)$.
In (4.2.29) and (4.2.30), $C$ depends on the a priori bound given by (4.2.11), while $E_{0}, C_{1}$ are uniform constants. Consequently, if we take $B_{1}$ large enough relatively to $\|f\|_{\mathscr{H} s+\zeta}, E_{0}, C_{1}$ and then $\epsilon \leq \delta^{2} \leq \delta_{0}^{2}$, with $\delta_{0}$ small enough, we deduce from (4.2.29) and (4.2.30) that (4.2.6) holds at rank $k+1$. The second estimate in (4.2.9) at rank $k+1$ follows, with for instance $B_{2}=2 B_{1}$, if $\delta_{0}$ is small enough. We are left with establishing the first estimate in (4.2.9) at rank $k+1$ and (4.2.10).

First let us bound $W_{k+1}-W_{k}$. By (4.2.25), for any $k^{\prime} \in\{1, \ldots, k\}$, any $(\omega, \epsilon) \in[1,2] \times\left[0, \delta^{2}\right]-G_{k^{\prime}}$, $\alpha \in \mathscr{A}-\mathscr{A}_{0}$, and any $2^{k^{\prime}} \leq\langle n(\alpha)\rangle<2^{k^{\prime}+1}$ (or any $(\omega, \epsilon) \in[1,2] \times\left[0, \delta^{2}\right]-G_{0}$ and $\alpha \in \mathscr{A}_{0}$ ), we have

$$
\begin{aligned}
& \left(\widetilde{L}_{\omega}+\epsilon V_{\mathrm{D}}\left(U_{k}, \omega, \epsilon\right)\right) \widetilde{\Pi}_{\alpha} W_{k+1}=\epsilon \widetilde{\Pi}_{\alpha} H_{k+1}\left(U_{k}, W_{k}\right), \\
& \left(\widetilde{L}_{\omega}+\epsilon V_{\mathrm{D}}\left(U_{k-1}, \omega, \epsilon\right)\right) \widetilde{\Pi}_{\alpha} W_{k}=\epsilon \widetilde{\Pi}_{\alpha} H_{k}\left(U_{k-1}, W_{k-1}\right),
\end{aligned}
$$

whence the equation

$$
\begin{align*}
\left(\widetilde{L}_{\omega}+\epsilon V_{\mathrm{D}}\left(U_{k}, \omega, \epsilon\right)\right) \widetilde{\Pi}_{\alpha}\left(W_{k+1}-W_{k}\right)=\epsilon \widetilde{\Pi}_{\alpha}[ & \left.V_{\mathrm{D}}\left(U_{k-1}, \omega, \epsilon\right)-V_{\mathrm{D}}\left(U_{k}, \omega, \epsilon\right)\right] W_{k} \\
& +\epsilon \widetilde{\Pi}_{\alpha}\left[H_{k+1}\left(U_{k}, W_{k}\right)-H_{k}\left(U_{k-1}, W_{k-1}\right)\right] \tag{4.2.31}
\end{align*}
$$

We make act $\mathscr{A}_{\alpha}\left(\omega ; U_{k}, \epsilon\right)^{-1}$ on both sides as in (4.2.25). Applying inequality (4.1.9) we get

$$
\begin{align*}
\left\|\tilde{\Pi}_{\alpha}\left(W_{k+1}-W_{k}\right)\right\|_{\mathscr{H} \sigma} \leq \frac{E_{0} \epsilon}{\delta}\left[\| \tilde{\Pi}_{\alpha}\left[V_{\mathrm{D}}( \right.\right. & \left.\left.U_{k-1}, \omega, \epsilon\right)-V_{\mathrm{D}}\left(U_{k}, \omega, \epsilon\right)\right] W_{k} \|_{\mathscr{H} \sigma+\zeta} \\
& \left.+\left\|\tilde{\Pi}_{\alpha}\left[H_{k+1}\left(U_{k}, W_{k}\right)-H_{k}\left(U_{k-1}, W_{k-1}\right)\right]\right\|_{\mathscr{H} \sigma+\zeta}\right] \tag{4.2.32}
\end{align*}
$$

This estimate holds outside $G_{k^{\prime}}$ when $k^{\prime} \neq 0, \alpha \in \mathscr{A}-\mathscr{A}_{0}, 2^{k^{\prime}} \leq\langle n(\alpha)\rangle<2^{k^{\prime}+1}$, and outside $G_{0}$ when $\alpha \in \mathscr{A}_{0}$. By (4.2.28), we may write

$$
\begin{align*}
& \left(W_{k+1}-W_{k}\right)(t, x, \omega, \epsilon)=\sum_{\alpha \in \mathscr{A}_{0}}\left(1-\psi_{0}\right) \widetilde{\Pi}_{\alpha}\left(W_{k+1}-W_{k}\right) \\
& \quad+\sum_{k^{\prime}=1}^{k} \sum_{\substack{\alpha \in \mathscr{A}-\mathscr{A}_{0} \\
2^{k^{\prime} \leq\langle n(\alpha)\rangle<2^{k^{\prime}+1}}}}\left(1-\psi_{k^{\prime}}\right) \widetilde{\Pi}_{\alpha}\left(W_{k+1}-W_{k}\right)+\sum_{\substack{\alpha \in \mathscr{A}-\mathscr{A}_{0} \\
2^{k+1} \leq\langle n(\alpha)\rangle<2^{k+2}}}\left(1-\psi_{k+1}\right) \widetilde{\Pi}_{\alpha} W_{k+1} . \tag{4.2.33}
\end{align*}
$$

The $\mathscr{H}^{\sigma}$ norm of the last term is bounded by $C_{2} 2^{-k(s-\sigma)}\left\|W_{k+1}\right\|_{\mathscr{H} s} \leq C_{2} B_{1} \frac{\epsilon}{\delta} 2^{-k(s-\sigma)}$ by (4.2.6), for some universal constant $C_{2}$. The $\mathscr{H}^{\sigma}$-norm of the $k^{\prime}$-sum in (4.2.33) may be estimated using (4.2.32), (4.2.20) and the bound

$$
\left\|\left(V_{\mathrm{D}}\left(U_{k-1}, \omega, \epsilon\right)-V_{\mathrm{D}}\left(U_{k}, \omega, \epsilon\right)\right) W_{k}\right\|_{\mathscr{H} \sigma+\zeta} \leq C\left\|U_{k}-U_{k-1}\right\|_{\mathscr{H} \sigma}\left\|W_{k}\right\|_{\mathscr{H}}
$$

which follows from (2.1.2), and where we used $s \geq \sigma+\zeta$. Using the induction hypothesis (4.2.9), (4.2.10), we get

$$
\begin{align*}
& \left\|W_{k+1}-W_{k}\right\|_{\mathscr{H} \sigma} \\
& \leq E_{0} \frac{\epsilon}{\delta}\left(C B_{1} \frac{\epsilon}{\delta} 2 B_{2} \frac{\epsilon}{\delta} 2^{-k \zeta}+3 C B_{2} \frac{\epsilon}{\delta} 2^{-k \zeta}+C 2^{-k \zeta}\left(B_{1}+B_{2}\right) \frac{\epsilon}{\delta}+(1+C \epsilon)\|f\|_{\mathscr{H} \sigma+2 \zeta} 2^{-k \zeta}\right) \\
& +C_{2} B_{1} \frac{\epsilon}{\delta} 2^{-k(s-\sigma)} \tag{4.2.34}
\end{align*}
$$

Since $s \geq \sigma+\zeta$, we may take $B_{1}$ large enough relatively to $E_{0},\|f\|_{\mathscr{H} s+\zeta}$, and $B_{2}$ large enough relatively to $C_{2}, B_{1}$, and $\epsilon / \delta \leq \delta \leq \delta_{0}$ small enough, so that (4.2.34) is smaller than $B_{2}(\epsilon / \delta) 2^{-(k+1) \zeta}$, whence (4.2.10) at rank $k+1$. Writing

$$
U_{k+1}-U_{k}=\left(\operatorname{Id}+\epsilon Q\left(U_{k}, \omega, \epsilon\right)\right)\left(W_{k+1}-W_{k}\right)+\epsilon\left(Q\left(U_{k}, \omega, \epsilon\right)-Q\left(U_{k-1}, \omega, \epsilon\right)\right) W_{k},
$$

we deduce from that the first inequality in (4.2.9) at rank $k+1$, for small enough $\epsilon$. This concludes the proof of the proposition.

Proof of Theorem 1.1.1. By (4.2.9), the series $\sum\left(U_{k}-U_{k-1}\right)$ converges in $\mathscr{H}^{\sigma}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}^{2}\right)$ and its sum $U$ satisfies $U \in \mathscr{H}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}^{2}\right)$ with

$$
\|U(\cdot, \omega, \epsilon)\|_{\mathscr{H} s}+\delta\left\|\partial_{\omega} U(\cdot, \omega, \epsilon)\right\|_{\mathscr{H} \text { s- }-2} \leq B_{2} \frac{\epsilon}{\delta}
$$

We have to check that $U$ gives a solution to our problem outside a set of parameters of small measure. Let

$$
(\omega, \epsilon) \in[1,2] \times\left[0, \delta^{2}\right]-\bigcup_{k^{\prime}=0}^{\infty} \mathcal{O}_{k^{\prime}}
$$

and $\delta \in] 0, \delta_{0}$ ]. Then (4.2.12) is satisfied for any $k$. We make $k \rightarrow+\infty$. Since we have uniform $\mathscr{H}^{s}$ bounds for $U_{k}, W_{k}$ and $\mathscr{H}^{\sigma}$ convergence for these quantities, the limit $U$ satisfies

$$
\left(\widetilde{L}_{\omega}+\epsilon V(U, \omega, \epsilon)\right) U=\epsilon R(U, \omega, \epsilon) U+\epsilon f
$$

which is (2.3.15). We have seen that this equation is equivalent to (2.3.14), which is, by Proposition 2.3.1, the same as (2.2.13). Since Proposition 2.2.4 shows that, up to a change of notation, this equation is equivalent to the formulation (2.2.6) of (1.1.3), we obtain a solution satisfying the requirements of Theorem 1.1.1. We still have to check that (1.1.5) holds with $\mathbb{O}=\bigcup_{k^{\prime}=0}^{\infty} \mathbb{O}_{k^{\prime}}$. According to (4.2.2), the set $\mathbb{O}_{0}$ is included in the set of those $(\omega, \epsilon)$ such that there are $(j, n)$ in a given finite subset of $\mathbb{Z}^{2}$ such that $\left|j \omega+|n|^{2}+\mu\right|<2 \delta$. The $\omega$-measure of this set is $O(\delta), \delta \rightarrow 0$ (Note that since $\mu \notin \mathbb{Z}_{-}$, we may always assume $j \neq 0$ ). For $k^{\prime}>0, \mathbb{O}_{k^{\prime}}$ is the union for $\alpha \in \mathscr{A}-\mathscr{A}_{0}$ with $2^{k^{\prime}-1} \leq\langle n(\alpha)\rangle<2^{k^{\prime}}$ and $\ell \in\left\{1, \ldots, D_{\alpha}\right\}$ of the set of those $(\omega, \epsilon)$ satisfying

$$
\left|\lambda_{\alpha}^{\ell}\left(\omega ; U_{k^{\prime}-1}, \epsilon\right)\right|<2 \delta 2^{-k^{\prime} \zeta}
$$

By (4.1.6), (4.1.7) the $\omega$-measure of each of these sets in bounded by $C\langle n(\alpha)\rangle^{-2} \delta 2^{-k^{\prime} \zeta} \leq C 2^{-\left(k^{\prime}+2\right) \zeta} \delta$. Since $D_{\alpha} \leq C_{1} 2^{k^{\prime}(\beta d+2)}$ by (1.2.4), (1.2.6), we obtain for the measure of the $\epsilon$-section of $\mathcal{O}$ the bound

$$
C \sum_{k^{\prime}=0}^{+\infty} 2^{-\left(k^{\prime}+2\right) \zeta+k^{\prime}(\beta d+2)+k^{\prime} d} \delta
$$

If we take $\zeta>(\beta+1) d+2$, we obtain the wanted $O(\delta)$ bound. This concludes the proof.

## Appendix

We gather here some elementary results used throughout the paper.
Lemma A.1. Let $s>\frac{d}{2}+1$. Then $\tilde{\mathscr{H}}^{s}\left(\mathbb{S}^{1} \times \mathbb{T}^{d} ; \mathbb{C}\right) \subset L^{\infty}$. Moreover, if $F$ is a smooth function on $\mathbb{S}^{1} \times \mathbb{T}^{d} \times \mathbb{C}$, satisfying $F(t, x, 0) \equiv 0$, there is some continuous function $\tau \rightarrow C(\tau)$ such that, for any $u \in \tilde{\mathscr{H}}^{s}$, we have $F(\cdot, u) \in \widetilde{\mathscr{H}}^{s}$ with the estimate

$$
\|F(\cdot, u)\|_{\mathscr{H}^{s}} \leq C\left(\|u\|_{L^{\infty}}\right)\|u\|_{\mathscr{H}^{s}}
$$

Proof. Let $\varphi \in C_{0}^{\infty}(] 0,+\infty[), \varphi \geq 0, \varphi \equiv 1$ on $[1,2]$ be such that $\sum_{\ell=-\infty}^{+\infty} \varphi\left(2^{-\ell} \lambda\right) \equiv 1$ for $\lambda \in \mathbb{R}_{+}^{*}$, and define $\psi(\lambda)=\sum_{-\infty}^{0} \varphi\left(2^{-\ell} \lambda\right)$. Consider, for $(j, n) \in \mathbb{Z} \times \mathbb{Z}^{d}$,

$$
\begin{align*}
& \Phi_{k}(j, n)=\varphi\left(2^{-2 k}\left(j^{2}+|n|^{4}\right)^{1 / 2}\right), k \geq 1, \\
& \Phi_{0}(j, n)=\psi\left(\left(j^{2}+|n|^{4}\right)^{1 / 2}\right) . \tag{A.1}
\end{align*}
$$

Define for $u \in \widetilde{\mathscr{H}}^{0}$ and $k \in \mathbb{N}$,

$$
\begin{equation*}
\Delta_{k} u=\sum_{j, n} \Phi_{k}(j, n) \hat{u}(j, n) \frac{e^{i(t j+k \cdot n)}}{(2 \pi)^{(d+1) / 2}}, \quad K_{k}(t, x)=\frac{1}{(2 \pi)^{d+1}} \sum_{j, n} \Phi_{k}(j, n) e^{i(t j+k \cdot n)} \tag{A.2}
\end{equation*}
$$

Then, for any $N \in \mathbb{N}$,

$$
\begin{equation*}
\left|K_{k}(t, x)\right| \leq C_{N} 2^{2 k(1+d / 2)}\left(1+2^{2 k}\left|e^{i t}-1\right|+2^{k}\left|e^{i x}-1\right|\right)^{-N} \tag{A.3}
\end{equation*}
$$

and $u \in \tilde{\mathscr{H}}^{s}$ if and only if $\left(2^{k s}\left\|\Delta_{k} u\right\|_{L^{2}}\right)_{k}$ is in $\ell^{2}$.
The first statement of the lemma follows from the inequality $\left\|\Delta_{k} u\right\|_{L^{\infty}} \leq C 2^{k(1+d / 2)}\left\|\Delta_{k} u\right\|_{L^{2}}$, which is a consequence of (A.3) (for the kernel corresponding to an enlarged $\Phi_{k}$ ). To get the second statement, we consider first the case of a function $F$ that does not depend on $(t, x)$. We set $S_{k}=$ $\sum_{k^{\prime} \leq k-1} \Delta_{k^{\prime}}$ when $k \geq 1, S_{0}=0$ and write

$$
F(u)=\sum_{k=0}^{+\infty}\left(F\left(S_{k+1} u\right)-F\left(S_{k} u\right)\right)=\sum_{k=0}^{+\infty} m_{k}(u) \Delta_{k} u
$$

where $m_{k}(u)=\int_{0}^{1} F^{\prime}\left(S_{k} u+\tau \Delta_{k} u\right) d \tau$. It follows from the definition of $S_{k}$ that this operator is given by a convolution kernel obeying the same estimates as in (A.3). Consequently, for any $(\alpha, \beta) \in \mathbb{N} \times \mathbb{N}^{d}$,

$$
\begin{equation*}
\left\|\partial_{t}^{\alpha} \partial_{x}^{\beta} m_{k}(u)\right\|_{L^{\infty}} \leq C 2^{2 k \alpha+k|\beta|} \tag{A.4}
\end{equation*}
$$

with a constant depending only on $\|u\|_{L^{\infty}}$. One writes for some $N_{0} \in \mathbb{N}$ to be chosen

$$
\begin{equation*}
\Delta_{j}[F(u)]=\sum_{k=0}^{j-1-N_{0}} \Delta_{j}\left[m_{k}(u) \Delta_{k} u\right]+\sum_{k=j-N_{0}}^{+\infty} \Delta_{j}\left[m_{k}(u) \Delta_{k} u\right] \tag{A.5}
\end{equation*}
$$

The $L^{2}$-norm of the second sum is bounded by $C c_{j} 2^{-j s}\|u\|_{\tilde{H} s}$ for some sequence $\left(c_{j}\right)_{j}$ in the unit ball of $\ell^{2}$, and some $C$ depending only on $\|u\|_{L^{\infty}}$. If $N_{0}$ is fixed large enough, because of the support properties of the Fourier transforms,

$$
\Delta_{j}\left[m_{k}(u) \Delta_{k} u\right]=\Delta_{j}\left[\left[\left(\mathrm{Id}-S_{j-N_{0}}\right) m_{k}(u)\right] \Delta_{k} u\right]
$$

when $k \leq j-1-N_{0}$. We estimate the $L^{2}$-norm of this quantity by

$$
\begin{equation*}
\left\|\left(\operatorname{Id}-S_{j-N_{0}}\right) m_{k}\right\|_{L^{\infty}}\left\|\Delta_{k} u\right\|_{L^{2}} \tag{A.6}
\end{equation*}
$$

and use that, for any $N$, we have $\left\|\left(\operatorname{Id}-S_{j-N_{0}}\right) m_{k}\right\|_{L^{\infty}} \leq C_{N} 2^{-4 j N}\left\|P^{N} m_{k}\right\|_{L^{\infty}}$ where $P=\partial_{t}^{2}+\Delta^{2}+1$. It follows from (A.4) that (A.6) is bounded from above by $C_{N} 2^{-4(j-k) N}\left\|\Delta_{k} u\right\|_{L^{2}}$, from which we deduce that the $L^{2}$-norm of the first sum in (A.5) is also smaller than $C 2^{-j s} c_{j}\|u\|_{\mathscr{H}_{s}}$. This concludes the proof for functions $F$ independent of $(t, x)$. In the general case, we note that since $u$ is bounded, we may always assume that $F$ is compactly supported, and we write

$$
F(t, x, u)=\frac{1}{2 \pi} \int_{\mathbb{R}} F_{1}(u, \theta) b(t, x, \theta) d \theta
$$

where $F_{1}(u, \theta)=e^{i u \theta}-1$ and $b(t, x, \theta)$ is the Fourier transform of $u \rightarrow F(t, x, u)$. Then it follows from the preceding proof that $F_{1}(u, \theta)$ is in $\tilde{\mathscr{H}}^{s}$ with a bound $\left\|F_{1}(u, \theta)\right\|_{\tilde{\mathscr{H}}^{s}} \leq C\langle\theta\rangle^{N(s)}$, for some exponent $N(s)$. Moreover, for any $N,\|b(\cdot, \theta)\|_{\mathscr{H}^{s}} \leq C_{N}\langle\theta\rangle^{-N}$. We get the conclusion by superposition.

Corollary A.2. Let $F: \mathbb{S}^{1} \times \mathbb{T}^{d} \times \mathbb{C} \rightarrow \mathbb{C}$ be a smooth function with $F(t, x, 0) \equiv 0$. Then $u \rightarrow F(\cdot, u)$ is a smooth map from $\tilde{\mathscr{H}}^{\sigma}$ to itself, for any $\sigma>\frac{d}{2}+1$.
Proof. We write

$$
F(t, x, u+h)-F(t, x, u)-\partial_{u} F(t, x, u) h=\int_{0}^{1} \int_{0}^{1}\left(D^{2} F\right)\left(t, x, u+\tau_{1} \tau_{2} h\right) \tau_{1} \cdot h^{2} d \tau_{1} d \tau_{2}
$$

and we apply the lemma to $D^{2} F(t, x, u)-D^{2} F(t, x, 0)$.
Lemma A.3. - Let $s>\frac{d}{2}+1$. If $u \in \tilde{\mathscr{H}}^{s}$ and $v \in \tilde{\mathscr{H}}^{\sigma^{\prime}}$ for some $\sigma^{\prime} \in[-s, s]$, then $u v \in \tilde{\mathscr{H}} \sigma^{\prime}$.

- For any $\sigma \in \mathbb{R}$ and $\sigma_{0}>\frac{d}{2}+1, \tilde{\mathscr{H}}^{\sigma} \cdot \tilde{\mathscr{H}}^{-\sigma} \subset \tilde{\mathscr{H}}^{-\max \left(\sigma, \sigma_{0}\right)}$.


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Received 16 Oct 2009. Revised 13 Jan 2010. Accepted 14 Sep 2010.
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APDE peer review and production are managed by EditFLOw ${ }^{\mathrm{TM}}$ from Mathematical Sciences Publishers.

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## Volume 4 No. 52011

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[^0]:    This work was partially supported by the ANR project Equa-disp.
    MSC2000: 35B10, 35Q55.
    Keywords: nonlinear Schrödinger equation, periodic solutions.

