# Sobolev periodic solutions of nonlinear wave equations in higher spatial dimensions 

Massimiliano Berti, ${ }^{*}$ Philippe Bolle ${ }^{\dagger}$


#### Abstract

We prove the existence of Cantor families of periodic solutions for nonlinear wave equations in higher spatial dimensions with periodic boundary conditions. We study both forced and autonomous PDEs. In the latter case our theorems generalize previous results of Bourgain to more general nonlinearities of class $C^{k}$ and assuming weaker non-resonance conditions. Our solutions have Sobolev regularity both in time and space. The proofs are based on a differentiable Nash-Moser iteration scheme, where it is sufficient to get estimates of interpolation type for the inverse linearized operators. Our approach works also in presence of very large "clusters of small divisors".


Keywords: Nonlinear Wave Equation, Infinite Dimensional Hamiltonian Systems, Nash-Moser Theorem, Periodic Solutions, Lyapunov-Schmidt reduction, small divisors ${ }^{1}$.
2000AMS subject classification: 35L05, 37K50, 37K55.

## 1 Introduction

The search for periodic solutions of nonlinear wave equations has a long standing tradition. The first pioneering results of Rabinowitz [23] and Brezis-Coron-Nirenberg [5] proved, by means of global variational methods, the existence of periodic solutions for 1-dimensional nonlinear wave equations, with a rational frequency. The reason for such a condition is that the other frequencies give rise to a "small divisors" problem type, due to complex resonance phenomena.

On the other hand, the existence of periodic and quasi periodic solutions in a neighborhood of an elliptic equilibrium, for positive measure sets of frequencies, was also considered. In this direction, the first results have been proved by Kuksin [16] and Wayne [25] for one dimensional, analytic, nonlinear wave equations. The main difficulty, namely the presence of arbitrarily "small divisors" in the expansion series of the solutions, is handled via KAM theory. The pioneering results in [16]-[25] were limited to Dirichlet boundary conditions because they required the eigenvalues of the Laplacian to be simple (the square roots of the eigenvalues are the normal modes frequencies of small oscillations of the string). In this case one can impose strong non-resonance conditions between the "tangential" and the "normal" frequencies of the expected KAM torus (the so-called "second order Melnikov" non-resonance conditions) to solve the linear homological equations which arise at each step of the KAM iteration, see also [22]-[18]-[17]. Such equations are linear PDEs with constant coefficients and can be solved by standard Fourier series. For periodic boundary conditions, where two consecutive eigenvalues are possibly equal, the second order Melnikov non-resonance conditions are violated.

In order to overcome such limitations, Craig and Wayne [11] introduced the Lyapunov-Schmidt decomposition method for PDEs and solved the small divisors problem, for periodic solutions, with an analytic Newton iteration scheme. Such an approach is particularly designed for dealing with resonant situations. On the other hand, the main difficulty of this strategy lies in the inversion of the linearized operators obtained at each step of the iteration, and in achieving suitable estimates for their inverse in

[^0]high (analytic) norms. Indeed these operators come from self-adjoint linear PDEs with non-constant coefficients and are small perturbations of a diagonal operator having arbitrarily small eigenvalues. To solve this problem Craig and Wayne imposed, for positive measure sets of parameters, lower bounds for the moduli of the eigenvalues (in particular, all the eigenvalues must be non-zero). These assumptions imply upper bounds for the operatorial $L^{2}$-norm of the inverse operators. Next, in order to get estimates on the inverse in analytic norms, Craig and Wayne developed a coupling technique inspired by the methods of Frölich-Spencer [14] in the Anderson localization theory, see also [10]. Key ingredients to achieve estimates for the inverse operators in high norms are the following assumptions on the unperturbed operator: ( $i$ ) "separation properties" between clusters of singular sites (that is between the Fourier indexes of the small divisors), (ii) properties of "well-localization" of the eigenfunctions with respect to the exponentials. These two facts, together with the analyticity of the functions, imply a very weak interaction between the singular sites. The second requirement $(i i)$ is free when working with periodic boundary conditions. On the other hand, since the first requirement $(i)$ can be obtained by imposing only the "first order Melnikov" non-resonance conditions, the Craig-Wayne approach works perfectly also in case of degenerate eigenvalues. In [11] the "clusters of small divisors" have a fixed bounded size: it is the case for one dimensional nonlinear wave and Schrödinger equations with periodic boundary conditions, see also [10].

The main difficulty in extending these results to PDEs in higher spatial dimensions is that the eigenvalues of the Laplacian can be highly degenerate, forming clusters of increasing size which tends to infinity.

This further problem has been first solved by Bourgain [6] for nonlinear wave equations in dimension $d \geq 2$ with periodic boundary conditions, extending the Craig-Wayne techniques. These results hold for analytic (polynomial) nonlinearities and prove the existence of periodic solutions having Gevrey regularity both in time and space. Suitable separations properties between the clusters of small divisors are imposed in [6] assuming a strong Diophantine-type condition, see remark 4.1. Then, using repeatedly the resolvent identity (see [7]), Bourgain proves a sub-exponentially fast decay for the off-diagonal terms of the inverse matrix. This estimate on the speed of decay gives an upper bound for the inverse matrix in Gevrey norm. For this step the high (Gevrey) regularity of the given functions is exploited. Similar results for the nonlinear Schrödinger equation have been proved in [7, Appendix 2], but, in that case, the separation conditions for the clusters of small divisors are more simply obtained.
The main result in [7] actually proves the existence also of quasi-periodic solutions in dimension $d=2$. See also [8] about the construction of quasi-periodic orbits for the nonlinear wave and Schrödinger equations in any spatial dimension.

In the present paper we prove the existence of periodic solutions for higher dimensional nonlinear wave equations for merely differentiable nonlinearities and under weaker non-resonance conditions than in [6]. We consider both forced and autonomous PDEs. In the forced case Theorem 1.1 is the first higher dimensional result, and extends [21], [13], [2], which are valid for 1-dimensional equations. In the autonomous case Theorem 1.2 generalizes the result of Bourgain [6]. Our solutions have the same Sobolev regularity both in time and space.

In order to prove our results we need all the power of the differentiable Nash-Moser theory. In particular, the key point of the iterative process lies in the "a-priori" bounds (15) for the divergence of the high Sobolev norms of the approximate solutions; we refer to [4] for further explanations and for a comparison with the approach of [11].

Concerning the linearized operators obtained at each step of the Nash-Moser iteration, it is sufficient to achieve just interpolation type estimates for their inverses, see the key property (P5). Our approach works also in presence of possibly very large clusters of small divisors: the "dyadic" condition (H1) (see section 3.2) is weaker than the corresponding ones in [7] (Lemma 7 of Appendix 2) and in [15]. Furthermore (H1) could also be considerably weakened, see remark 3.1, even though not completely eliminated, see the discussion below and remark 3.2. A point of interest is that, in presence of possibly very large clusters of small divisors, the use of Sobolev norms, instead of analytic or Gevrey ones, used in [11]-[6], makes the estimates easier. The most intuitive reason is that a lower bound for the moduli of the eigenvalues yields immediately a $L^{2}$-bound for the inverse matrix, and the Sobolev norms are closer to the $L^{2}$-norm than the

Gevrey or analytic norms, see lemma 3.1 and related comments. Clearly, working with functions having a mere Sobolev regularity, has the drawback of a slow (polynomial) decay off the diagonal of the matrix elements of these operators. This makes the interactions between their singular clusters rather strong. Nevertheless a polynomial decay of large enough order (connected to some smoothness assumption) is sufficient.

On the other hand, we underline that it does not seem sufficient to have only lower bounds for the moduli of all the eigenvalues without some separation properties between the singular clusters. This information, by itself, would give a too weak estimate for the norm of the inverse matrix, and the NashMoser scheme would not converge, see remark 3.2. This is also related to a famous counter-example of Lojasiewicz-Zehnder [19] concerning the optimal conditions in an abstract Nash-Moser implicit function theorem.

Since the aim of the present paper is to focus on the solution of the small divisors problem in presence of large clusters and with differentiable nonlinearities, we have considered model cases in which the bifurcation equation arising with the Lyapunov-Schmidt reduction, or is not present (as for Theorem 1.1) or it is rather simply solved (as for Theorem 1.2).

Before concluding this introduction, we mention that the KAM approach has been extended by Chierchia-You [9] to prove the existence of quasi-periodic solutions for one dimensional wave and Schrödinger equations in case of periodic boundary conditions (in particular, their theorem covers the result of [11]), and by Eliasson-Kuksin [12] for higher dimensional nonlinear Schrödinger equations. We also mention that Gentile-Procesi [15] have recently obtained the existence of periodic solutions for higher dimensional nonlinear Schrödinger equations by the Lindstedt series method. We remark that in all the previous results the nonlinearities are required to be analytic and the solutions are analytic in time.
Acknowledgments: The authors thank P. Baldi, L. Biasco and M.Procesi for interesting comments.

### 1.1 Main results

Let us consider $d$-dimensional nonlinear wave equations with periodic boundary conditions of the form

$$
\left\{\begin{array}{l}
u_{t t}-\Delta u+m u=\varepsilon F(\omega t, x, u)  \tag{1}\\
u(t, x)=u(t, x+2 \pi k), \quad \forall k \in \mathbf{Z}^{d}
\end{array}\right.
$$

where the forcing term $F(\omega t, x, u)$ is $2 \pi / \omega$-periodic in time ${ }^{2}, 2 \pi$-periodic in each spatial variable $x_{i}$, $i=1, \ldots, d, m \in \mathbf{R}$ and $\varepsilon>0$ is a small parameter.

We consider the non-resonant case when

$$
\begin{equation*}
\left|\omega^{2} l^{2}-\lambda_{j}\right| \geq \frac{3 \gamma}{\max \left(1,|l|^{3 / 2}\right)}, \quad \gamma \in(0,1), \quad \lambda_{j}:=|j|^{2}+m, \quad \forall(l, j) \in \mathbf{Z} \times \mathbf{Z}^{d} . \tag{2}
\end{equation*}
$$

Note that in (2), the exponent $3 / 2$ is fixed for simplicity of exposition and could be replaced by any $\alpha>1$. For all $j$ such that $\lambda_{j} \geq 0$ define $\omega_{j}:=\sqrt{\lambda_{j}}$. Assumption (2) means that the forcing frequency $\omega$ does not enter in resonance with the normal mode frequencies $\omega_{j}$ of oscillations of the membrane.

By standard arguments, (2) is satisfied for all $\omega$ in $\left[\bar{\omega}_{1}, \bar{\omega}_{2}\right]$ but a subset of measure $O(\gamma)$.
Concerning regularity we only assume that $F \in C^{k}\left(\mathbf{T} \times \mathbf{T}^{d} \times \mathbf{R} ; \mathbf{R}\right)$ for some $k$ large enough.
If $F(t, x, 0) \not \equiv 0$ then $u=0$ is not a solution of (1) for $\varepsilon \neq 0$.

- Question: do there exist periodic solutions of (1) for positive measure sets of $(\varepsilon, \omega)$ ?

Normalizing the period, we look for $2 \pi$-periodic in time solutions of

$$
\begin{equation*}
\omega^{2} u_{t t}-\Delta u+m u=\varepsilon F(t, x, u) \tag{3}
\end{equation*}
$$

[^1](where $F$ is $2 \pi$-periodic in $t$ ) in the real Sobolev space ${ }^{3}$
\[

$$
\begin{aligned}
H^{s}:=H^{s}\left(\mathbf{T} \times \mathbf{T}^{d} ; \mathbf{R}\right):= & \left\{u(t, x):=\sum_{(l, j) \in \mathbf{Z} \times \mathbf{Z}^{d}} u_{l, j} e^{\mathrm{i}(l t+j \cdot x)}, u_{l, j}^{*}=u_{-l,-j}\right. \\
& \left.\left.\left|\|u\|_{s}^{2}:=\sum_{(l, j) \in \mathbf{Z} \times \mathbf{Z}^{d}}\right| u_{l, j}\right|^{2} \max \left(\left(l^{2}+|j|^{2}\right)^{s}, 1\right)<+\infty\right\}
\end{aligned}
$$
\]

for some $s>(d+1) / 2$.
We recall that for $s>(d+1) / 2$ we have the continuous inclusion $H^{s}\left(\mathbf{T}^{d+1}\right) \hookrightarrow L^{\infty}\left(\mathbf{T}^{d+1}\right)$ and $H^{s}$ is a Banach algebra with respect to the multiplication of functions.

Theorem 1.1 Fix $0<\bar{\omega}_{1}<\bar{\omega}_{2}$. There is $s:=s(d), k:=k(d) \in \mathbf{N}$, such that: $\forall F \in C^{k}\left(\mathbf{T} \times \mathbf{T}^{d} \times \mathbf{R}\right), \forall \gamma \in(0,1)$, there exist $\varepsilon_{0}:=\varepsilon_{0}(\gamma), K, C>0$ (independent of $\gamma$ ), a map

$$
\widetilde{u} \in C^{1}\left(\left[0, \varepsilon_{0}\right] \times\left[\bar{\omega}_{1}, \bar{\omega}_{2}\right] ; H^{s}\right) \quad \text { with } \quad\|\widetilde{u}(\varepsilon, \omega)\|_{s} \leq K \gamma^{-1} \varepsilon,\left\|\left(D_{\varepsilon, \omega} \widetilde{u}\right)(\varepsilon, \omega)\right\|_{s} \leq K \gamma^{-1}
$$

and a Cantor like set $A_{\infty} \subset\left[0, \varepsilon_{0}\right] \times\left[\bar{\omega}_{1}, \bar{\omega}_{2}\right]$, of Lebesgue measure

$$
\begin{equation*}
\left|A_{\infty}\right| \geq \varepsilon_{0}\left(\left|\bar{\omega}_{2}-\bar{\omega}_{1}\right|-C \gamma\right) \tag{4}
\end{equation*}
$$

such that, $\forall(\varepsilon, \omega) \in A_{\infty}, \widetilde{u}(\varepsilon, \omega)$ is a solution of (3).
Note that, as the freely chosen small parameter $\gamma$ tends to 0 , the constant $\varepsilon_{0}(\gamma)$ tends to 0 , and the bounds on $\widetilde{u}$ get worse, but, as a counterpart, the ratio $\left|A_{\infty}\right| / \varepsilon_{0}\left(\bar{\omega}_{2}-\bar{\omega}_{1}\right)$ tends to 1 , i.e. the set $A_{\infty}$ has asymptotically full measure.

The conditions defining $A_{\infty}$ are (2) and (7) (which are independent of $\varepsilon$ ), plus infinitely many others, which depend on the nonlinearity, and are required to get the invertibility of the linearized operators obtained at each step of the Nash-Moser iteration (see Theorem 2.1).

Remark 1.1 The non-resonance condition (2) implies $m \neq 0$. Note that if $m<0$ the equilibrium $u=0$ is not completely elliptic. In the case $m=0$, under some additional assumptions on $F$, a result similar to Theorem 1.1 holds assuming condition (2) only for all $(l, j) \neq 0$. Then we perform a Lyapunov-Schmidt reduction according to the decomposition $H^{s}=\mathbf{R} \oplus H_{0}^{s}$ where $H_{0}^{s}$ denote the Sobolev functions with zero mean value.

We develop in detail all the computations to prove Theorem 1.1. The same techniques can be used to prove the existence of Cantor families of small amplitude periodic solutions for autonomous $d$-dimensional nonlinear wave equations of the form

$$
\left\{\begin{array}{l}
u_{t t}-\Delta u+m u=a u^{p}+r(x, u)  \tag{5}\\
u(t, x)=u(t, x+2 \pi k), \quad \forall k \in \mathbf{Z}^{d}
\end{array}\right.
$$

with $p \geq 3$ odd integer, $a \in \mathbf{R}, a \neq 0$, and

$$
\begin{equation*}
r(x, u) \in C^{k}\left(\mathbf{T}^{d} \times \mathbf{R}\right), k \geq p, \quad r(x, 0)=\ldots=\partial_{u}^{p} r(x, 0)=0 \quad \text { and } \quad r(-x, u)=r(x, u) \tag{6}
\end{equation*}
$$

The bounded solutions of the linearized equation $u_{t t}-\Delta u+m u=0$ which are even in time and in $x$ are

$$
u=\sum_{j \in \mathbf{Z}^{d},|j|^{2}+m \geq 0} \cos \left(\omega_{j} t\right) A_{j} \cos (j \cdot x), \quad A_{j} \in \mathbf{R}, \quad \omega_{j}:=\sqrt{|j|^{2}+m}
$$

Fixed $j_{0}$, we aim to prove the existence of small amplitude periodic solutions of the nonlinear equation (5) with frequencies close to $\omega_{j_{0}}$. Assuming that $m$ is irrational, the normal mode frequencies $\omega_{j}$ commensurable with $\omega_{j_{0}}$ satisfy $|j|=\left|j_{0}\right|$. We shall assume the stronger hypothesis that $m$ is Diophantine, to have a quantitative non-resonance condition like

$$
\left|\omega_{j}^{2}-l^{2} \omega_{j_{0}}^{2}\right| \geq \frac{\gamma}{(|l|+1)^{\tau}}, \quad \forall(|j|,|l|) \neq\left(\left|j_{0}\right|, 1\right)
$$

[^2]similar to (2), see [6].
Rescaling the amplitude $u \mapsto \delta u, \delta>0$ and normalizing the period, we look for solutions of
$$
\omega^{2} u_{t t}-\Delta u+m u=\varepsilon g(\delta, x, u), \quad \varepsilon:=\delta^{p-1}, \quad g(\delta, x, u):=a u^{p}+\delta^{-p} r(x, \delta u)
$$
in the subspace $H_{\text {even }}^{s}:=\left\{u \in H^{s} \mid u(-t,-x)=u(t, x)\right\}$. The regularity property of the composition operator induced by $g$ on $H_{\text {even }}^{s}$ are proved as in [4].

Theorem 1.2 Let $p \geq 3$ be an odd integer and assume (6). Suppose $m>0$ is Diophantine. Fix $j_{0} \in \mathbf{Z}^{d}$. There is $s:=s(d), k:=k(d) \in \mathbf{N}$ such that, for all $r \in C^{k}, \forall \gamma \in(0,1)$, there exist $\delta_{0}>0, \bar{A}_{j_{0}} \in \mathbf{R}, a$ curve

$$
u \in C^{1}\left(\left[0, \delta_{0}\right] ; H_{\text {even }}^{s}\right) \quad \text { with } \quad\left\|u(\delta)-\delta \bar{A}_{j_{0}} \cos (t) \cos \left(j_{0} \cdot x\right)\right\|_{s}=O\left(\delta^{2}\right),
$$

and a Cantor set $\mathcal{C} \subset\left[0, \delta_{0}\right]$ of asymptotically full measure such that, $\forall \delta \in \mathcal{C}, u(\delta)$ is a solution of $\omega^{2} u_{t t}-\Delta u+m u=u^{p}+r(x, u)$ with $\omega^{2}=\omega_{j_{0}}^{2}-\operatorname{sign}(a) \delta^{p-1}$.

Theorem 1.2 generalizes Bourgain's result in [6] to the case of differentiable nonlinearities, when the leading nonlinear term is $a u^{p}, p \geq 3$ and for irrational frequencies satisfying only (7). Furthermore the existence result of theorem 1.2 holds for a Cantor set $\mathcal{C}$ of asymptotically full measure. This requires to write the dependence with respect to $\gamma$ in the separation arguments of section 4 (used also in Theorem 1.1).

As said in the introduction the most difficult step in the proof of both Theorems 1.1-1.2 is to get estimates for the inverse linearized operators which arise at each step of the Nash-Moser iteration. For this task some "separation properties" for the "singular sites" seem required. For one dimensional wave equations, by a simple argument in [11]-[10], it is sufficient to assume that $\omega^{2}$ is Diophantine, namely

$$
\begin{equation*}
\left|\omega^{2} q-p\right| \geq \frac{\gamma}{\max \left(1,|p|^{3 / 2}\right)}, \quad \forall(q, p) \in \mathbf{Z}^{2} \backslash\{(0,0)\} \tag{7}
\end{equation*}
$$

with $\gamma \in(0,1)$. On the other hand, such separation properties are far from obvious in higher spatial dimension. They have been obtained by Bourgain in [6] under strong non-resonance conditions of Diophantine type for $\omega^{2}$, see remark 4.1. In section 4 we shall obtain the same separation properties assuming only (7). Also in (7), the exponent $3 / 2$ is fixed only for simplicity and could be replaced by any $\alpha>1$.

A final comment regarding the boundary conditions. The case of Dirichlet boundary conditions on a rectangle (with some oddness assumption on the nonlinearity) works similarly to the setting of the present case. The eigenfunctions of the unperturbed operator are still linear combinations of exponentials and the high order regularity of $u$ has a straightforward translation into the behavior of its Fourier coefficients, i.e. of its components in the orthonormal basis of the eigenfunctions. On the contrary, in the case of a general bounded domain $\Omega \subset \mathbf{R}^{d}$, the eigenfunctions of the Laplacian do not possess such a good property, even if they form an orthonormal basis of $H^{1}(\Omega)$. Therefore the existence of periodic solutions is this case is completely open (but see [1] in the case of an integral nonlinearity which does not mix the spatial modes).
Notations: $\mathcal{N}(A, \eta)$ denotes the $\eta$-neighborhood of a subset $A$ of a normed space; $z^{*}$ is the complex conjugate of $z \in \mathbf{C}$; the symbol $[x] \in \mathbf{N}$ denotes the integer part of $x \in \mathbf{R}$. We denote by $\mathcal{L}\left(H_{A}, H_{B}\right)$ the set of continuous linear operators from $H_{A}$ to $H_{B} . \mathrm{d}(A ; B):=\inf \{|a-b|, a \in A, b \in B\}$ is the distance from the set $A$ to the set $B$.

## 2 The Nash-Moser scheme

Consider the orthogonal splitting

$$
H^{s}=W^{\left(N_{n}\right)} \oplus W^{\left(N_{n}\right) \perp}
$$

where

$$
W^{\left(N_{n}\right)}=\left\{u \in H^{s} \mid u=\sum_{|(l, j)| \leq N_{n}} u_{l, j} e^{\mathrm{i}(l t+j \cdot x)}\right\}
$$

$$
W^{\left(N_{n}\right) \perp}=\left\{u \in H^{s} \mid u=\sum_{|(l, j)|>N_{n}} u_{l, j} e^{\mathrm{i}(l t+j \cdot x)}\right\}
$$

with ${ }^{4}$

$$
\begin{equation*}
N_{n}:=\left[e^{\lambda 2^{n}}\right], \quad \lambda:=\ln N_{0} \tag{8}
\end{equation*}
$$

In the proof of Theorem 2.1 we shall take $\lambda$, i.e. $N_{0} \in \mathbf{N}$, large enough. We denote by

$$
P_{N_{n}}: H^{s} \rightarrow W^{\left(N_{n}\right)} \quad \text { and } \quad P_{N_{n}}^{\perp}: H^{s} \rightarrow W^{\left(N_{n}\right) \perp}
$$

the orthogonal projectors onto $W^{\left(N_{n}\right)}$ and $W^{\left(N_{n}\right) \perp}$.
The convergence of the Nash-Moser scheme is based on properties $(P 1),(P 2),(P 3),(P 4),(P 5)$ below.
The first three properties are standard for the composition operator $f: H^{s} \rightarrow H^{s}$ defined by

$$
f(u)(t, x):=F(t, x, u(t, x))
$$

where $F \in C^{k}\left(\mathbf{T} \times \mathbf{T}^{d} \times \mathbf{R} ; \mathbf{R}\right)$ with $k \geq s+2$ and $s>(d+1) / 2$.

- (P1) (Regularity) $f \in C^{2}\left(H^{s} ; H^{s}\right)$ and $D^{2} f$ is bounded on $\left\{\|u\|_{s} \leq 1\right\}$.
- (P2) (Tame) $\forall s \leq s^{\prime} \leq k, \forall u \in H^{s^{\prime}}$ such that $\|u\|_{s} \leq 1,\|f(u)\|_{s^{\prime}} \leq C\left(s^{\prime}\right)\left(1+\|u\|_{s^{\prime}}\right)$.
- (P3) (Taylor Tame) $\forall s \leq s^{\prime} \leq k-2, \forall u \in H^{s^{\prime}}$ such that $\|u\|_{s} \leq 1, \forall h \in H^{s^{\prime}}$,

$$
\|f(u+h)-f(u)-D f(u) h\|_{s^{\prime}} \leq C\left(s^{\prime}\right)\left(\|u\|_{s^{\prime}}\|h\|_{s}^{2}+\|h\|_{s}\|h\|_{s^{\prime}}\right)
$$

where $[D f(u) h](t, x)=\left(\partial_{u} F\right)(t, x, u) h(t, x)$. In particular, for $s^{\prime}=s$,

$$
\begin{equation*}
\|f(u+h)-f(u)-D f(u) h\|_{s} \leq C\|h\|_{s}^{2} . \tag{9}
\end{equation*}
$$

We refer to [20] for the proof of $(P 2)$, see also [24]. Properties $(P 1)$ and $(P 3)$ are obtained similarly. Furthermore, by the definitions of the spaces $H^{s}$ and the projectors $P_{N}$, we have:

- (P4) (Smoothing) $\forall N \in \mathbf{N} \backslash\{0\}$,

$$
\begin{array}{ll}
\left\|P_{N} u\right\|_{s+r} \leq N^{r}\|u\|_{s}, \quad \forall u \in H^{s} \\
\left\|P_{N}^{\perp} u\right\|_{s} \leq N^{-r}\|u\|_{s+r}, \quad \forall u \in H^{s+r}
\end{array}
$$

The key property ( $P 5$ ), proved in section 3, is an invertibility property for the linearized operator

$$
\begin{equation*}
\mathcal{L}_{N}(\varepsilon, \omega, u(\varepsilon, \omega))[h]:=L_{\omega} h-\varepsilon P_{N} D f(u(\varepsilon, \omega)) h, \quad \forall h \in W^{(N)} \tag{10}
\end{equation*}
$$

where

$$
L_{\omega}:=\omega^{2} \partial_{t t}-\Delta+m
$$

and $u \in C^{1}\left(\left[0, \varepsilon_{0}\right] \times\left[\bar{\omega}_{1}, \bar{\omega}_{2}\right], W^{(N)}\right)$ with $\|u\|_{C^{1}(s)}:=\sup _{\left[0, \varepsilon_{0}\right] \times\left[\bar{\omega}_{1}, \bar{\omega}_{2}\right]}\|u(\varepsilon, \omega)\|_{s}+\left\|D_{\varepsilon, \omega} u(\varepsilon, \omega)\right\|_{s}$.
Let

$$
G:=\left\{(\varepsilon, \omega) \in\left[0, \varepsilon_{0}\right] \times\left[\bar{\omega}_{1}, \bar{\omega}_{2}\right] \mid \omega \text { satisfies (2) and (7) }\right\} .
$$

- (P5) (Invertibility of $\left.\mathcal{L}_{N}\right) \exists \mu:=\mu(d), s:=s(d), \bar{C}, c, K$, such that $\forall \gamma \in(0,1), \forall s^{\prime} \geq s, \forall C>0$, there exist $\varepsilon_{0}:=\varepsilon_{0}\left(\gamma, s^{\prime}, C\right)>0$ with the following property:
$\forall N, \forall u \in C^{1}\left(\left[0, \varepsilon_{0}\right] \times\left[\bar{\omega}_{1}, \bar{\omega}_{2}\right], W^{(N)}\right)$ with $\|u\|_{C^{1}(s)} \leq C \gamma^{-1}$, there is a set $G_{N}(u) \subset\left[0, \varepsilon_{0}\right] \times\left[\bar{\omega}_{1}, \bar{\omega}_{2}\right]$ such that, if $(\varepsilon, \omega) \in \mathcal{N}\left(G_{N}(u), \gamma N^{-\sigma}\right), \sigma:=\mu+3$, then $\mathcal{L}_{N}(\varepsilon, \omega, u(\varepsilon, \omega))$ is invertible and

$$
\begin{align*}
\left\|\mathcal{L}_{N}^{-1}(\varepsilon, \omega, u)[h]\right\|_{s^{\prime}} & \leq \frac{K\left(s^{\prime}\right)}{\gamma} N^{\mu}\left(\|h\|_{s^{\prime}}+\|u\|_{s^{\prime}}\|h\|_{s}\right), \quad \forall h \in W^{(N)}  \tag{11}\\
& \left\|\mathcal{L}_{N}^{-1}(\varepsilon, \omega, u)[h]\right\|_{s} \leq \frac{K}{\gamma} N^{\mu}\|h\|_{s} . \tag{12}
\end{align*}
$$

Moreover, if $N \leq c\left(\gamma \varepsilon_{0}^{-1}\right)^{2 / 3}$ then $G_{N}(u)=G$, and if $\left\|u_{1}-u_{2}\right\|_{s} \leq N^{-\sigma}$, then

$$
\begin{equation*}
\left|G_{N^{\prime}}^{c}\left(u_{2}\right) \backslash G_{N}^{c}\left(u_{1}\right)\right| \leq \bar{C} \frac{\gamma \varepsilon_{0}}{N}, \quad \forall N^{\prime} \geq N \tag{13}
\end{equation*}
$$

[^3]Theorem 2.1 (Nash-Moser) Let $0<\varepsilon \leq \varepsilon_{0}(\gamma)$ be small enough (possibly depending on d, $F, \bar{\omega}_{1}$, $\bar{\omega}_{2}$ ). There exists a decreasing sequence of subsets of "non-resonant" parameters

$$
\ldots \subseteq A_{n} \subseteq A_{n-1} \subseteq \ldots \subseteq A_{1} \subseteq A_{0}:=G \subset A:=\left[0, \varepsilon_{0}\right] \times\left[\bar{\omega}_{1}, \bar{\omega}_{2}\right]
$$

and a sequence of functions $\widetilde{u}_{n}(\varepsilon, \omega) \in C^{1}\left(A, W^{\left(N_{n}\right)}\right)$ satisfying $\left\|\widetilde{u}_{n}\right\|_{s} \leq K \varepsilon \gamma^{-1},\left\|D_{\varepsilon, \omega} \widetilde{u}_{n}\right\|_{s} \leq K \gamma^{-1}$, such that, $\forall n$, if $(\varepsilon, \omega) \in \mathcal{N}\left(A_{n}, \gamma N_{n}^{-\sigma} / 2\right)$ then $\widetilde{u}_{n}(\varepsilon, \omega)$ is a solution of
$\left(P_{N_{n}}\right)$

$$
L_{\omega} u-\varepsilon P_{N_{n}} f(u)=0 .
$$

On the Cantor like set $A_{\infty}:=\cap_{n \geq 0} A_{n}$, which satisfies the measure estimate (4), the sequence ( $\widetilde{u}_{n}(\varepsilon, \omega)$ ) converges in $H^{s}$ to a solution $\widetilde{u}(\varepsilon, \omega)$ of equation (3) satisfying $\|\widetilde{u}\|_{s} \leq K \varepsilon \gamma^{-1},\left\|D_{\varepsilon, \omega} \widetilde{u}\right\|_{s} \leq K \gamma^{-1}$. The convergence is uniform in $(\varepsilon, \omega) \in A_{\infty}$.

Proof. In the proof $K, K^{\prime}, K_{i}$, denote constants depending on $d, F, N_{0}, \bar{\omega}_{1}, \bar{\omega}_{2}$ at most.
First step: initialization. Assume $\sigma>7 / 2$. If $(\varepsilon, \omega) \in \mathcal{N}\left(G, \gamma N_{0}^{-\sigma}\right)$ then, by $(2), \forall|(l, j)| \leq N_{0}$, $\left|\omega^{2} l^{2}-\lambda_{j}\right| \geq \gamma / N_{0}^{3 / 2}$, for $N_{0}$ large so that $N_{0}^{\sigma-(7 / 2)}>\bar{\omega}_{2}$. Hence $\left\|L_{\omega}^{-1} h\right\|_{s} \leq N_{0}^{3 / 2} \gamma^{-1}\|h\|_{s}, \forall h \in W^{\left(N_{0}\right)}$.

By the contraction mapping theorem, using ( $P 1$ ), for $\varepsilon \gamma^{-1}$ small, there exists a unique solution $u_{0}:=u_{0}(\varepsilon, \omega)$ of equation $\left(P_{N_{0}}\right)$ satisfying $\left\|u_{0}\right\|_{s} \leq K_{0} \varepsilon \gamma^{-1}$. Furthermore, by the implicit function theorem, $(\varepsilon, \omega) \mapsto u_{0}(\varepsilon, \omega)$ is in $C^{1}\left(\mathcal{N}\left(G, \gamma N_{0}^{-\sigma}\right), H^{s}\right)$ and $\left\|D_{\varepsilon, \omega} u_{0}\right\|_{s} \leq K_{0} \gamma^{-1}$.

We define $\widetilde{u}_{0}:=\psi_{0} u_{0}$ where $\psi_{0}$ is a $C^{\infty}$ cut-off function defined on $A$ that takes the values 1 on $\mathcal{N}\left(G, \gamma N_{0}^{-\sigma} / 2\right)$ and 0 outside $\mathcal{N}\left(G, \gamma N_{0}^{-\sigma}\right)$, and $\left|D \psi_{0}\right| \leq C N_{0}^{\sigma} \gamma^{-1}$. By the previous estimates, $\left\|\widetilde{u}_{0}\right\|_{s} \leq K \varepsilon \gamma^{-1},\left\|D_{\varepsilon, \omega} \widetilde{u}_{0}\right\|_{s} \leq K \gamma^{-1}$, taking $\varepsilon C N_{0}^{\sigma} \gamma^{-1}<1$.
Second step: iteration. Suppose we have already defined $\widetilde{u}_{n} \in C^{1}\left(A, W^{\left(N_{n}\right)}\right)$ satisfying the properties stated in the proposition and, in addition, $\forall 1 \leq k \leq n$,

$$
\begin{gather*}
\left\|\widetilde{u}_{k}-\widetilde{u}_{k-1}\right\|_{s} \leq K \varepsilon \gamma^{-1} N_{k}^{-\sigma-1}, \quad\left\|D_{\varepsilon, \omega} \widetilde{u}_{k}-D_{\varepsilon, \omega} \widetilde{u}_{k-1}\right\|_{s} \leq K \gamma^{-1} N_{k}^{-1}  \tag{14}\\
B_{k} \leq\left(1+N_{k}^{\mu}\right) B_{k-1}, \quad \quad B_{k}^{\prime} \leq B_{k-1}^{\prime}+K \gamma^{-1} N_{k}^{\mu}\left(N_{k}^{\sigma} B_{k}+\varepsilon B_{k-1}^{\prime}\right) \tag{15}
\end{gather*}
$$

where

$$
B_{n}:=1+\left\|\widetilde{u}_{n}\right\|_{s+\beta}, \quad B_{n}^{\prime}:=1+\left\|D_{\varepsilon, \omega} \widetilde{u}_{n}\right\|_{s+\beta}
$$

and

$$
\begin{equation*}
\beta:=2(\sigma+1+3 \mu), \quad \sigma=\mu+3 \tag{16}
\end{equation*}
$$

By the first inequality in (15) we get

$$
\begin{equation*}
B_{n} \leq B_{0} \Pi_{i=1}^{n}\left(1+N_{i}^{\mu}\right) \leq K N_{n+1}^{\mu} \tag{17}
\end{equation*}
$$

and so, by the second inequality,

$$
\begin{equation*}
B_{n}^{\prime} \leq K \gamma^{-1} N_{n+1}^{2 \mu+\sigma} \tag{18}
\end{equation*}
$$

For $h \in W^{\left(N_{n+1}\right)}$ we write

$$
\begin{aligned}
L_{\omega}\left(\widetilde{u}_{n}+h\right)-\varepsilon P_{N_{n+1}} f\left(\widetilde{u}_{n}+h\right) & =\left[L_{\omega} \widetilde{u}_{n}-\varepsilon P_{N_{n+1}} f\left(\widetilde{u}_{n}\right)\right]+\left[L_{\omega} h-\varepsilon P_{N_{n+1}} D f\left(\widetilde{u}_{n}\right) h\right]+R_{n}(h) \\
& =r_{n}+\mathcal{L}_{N_{n+1}}\left(\varepsilon, \omega, \widetilde{u}_{n}\right) h+R_{n}(h)
\end{aligned}
$$

where

$$
r_{n}:=L_{\omega} \widetilde{u}_{n}-\varepsilon P_{N_{n+1}} f\left(\widetilde{u}_{n}\right) \quad \text { and } \quad R_{n}(h):=-\varepsilon P_{N_{n+1}}\left(f\left(\widetilde{u}_{n}+h\right)-f\left(\widetilde{u}_{n}\right)-D f\left(\widetilde{u}_{n}\right) h\right)
$$

If $(\varepsilon, \omega) \in \mathcal{N}\left(A_{n} ; \gamma N_{n}^{-\sigma} / 2\right)$ then $\widetilde{u}_{n}$ solves equation $\left(P_{N_{n}}\right)$ and so

$$
\begin{equation*}
r_{n}:=L_{\omega} \widetilde{u}_{n}-\varepsilon P_{N_{n+1}} f\left(\widetilde{u}_{n}\right)=-\varepsilon P_{N_{n}}^{\perp} P_{N_{n+1}} f\left(\widetilde{u}_{n}\right) \in W^{\left(N_{n}\right) \perp} \cap W^{\left(N_{n+1}\right)} \tag{19}
\end{equation*}
$$

Inversion of $\mathcal{L}_{N_{n+1}}\left(\varepsilon, \omega, \widetilde{u}_{n}\right)$. By property (P5), defining

$$
\begin{equation*}
A_{n+1}:=A_{n} \cap G_{N_{n+1}}\left(\widetilde{u}_{n}\right) \tag{20}
\end{equation*}
$$

$\mathcal{L}_{N_{n+1}}\left(\varepsilon, \omega, \widetilde{u}_{n}\right)$ is invertible for $(\varepsilon, \omega) \in \mathcal{N}\left(A_{n+1}, \gamma N_{n+1}^{-\sigma}\right)$ and

$$
\begin{equation*}
\left\|\mathcal{L}_{N_{n+1}}^{-1}\left(\varepsilon, \omega, \widetilde{u}_{n}\right) h\right\|_{s} \stackrel{(12)}{\leq} \frac{K}{\gamma} N_{n+1}^{\mu}\|h\|_{s}, \quad \forall h \in W^{\left(N_{n+1}\right)} \tag{21}
\end{equation*}
$$

Set

$$
\mathcal{G}_{n+1}: W^{\left(N_{n+1}\right)} \rightarrow W^{\left(N_{n+1}\right)}, \quad \mathcal{G}_{n+1}(h):=-\mathcal{L}_{N_{n+1}}^{-1}\left(\varepsilon, \omega, \widetilde{u}_{n}\right)\left(r_{n}+R_{n}(h)\right)
$$

so that equation $\left(P_{N_{n+1}}\right)$ is equivalent to the the fixed point problem $h=\mathcal{G}_{n+1}(h)$.
Lemma 2.1 (Contraction) $\exists K_{1}>0$ such that, for $\varepsilon \gamma^{-1}$ small enough, $\forall(\varepsilon, \omega) \in \mathcal{N}\left(A_{n+1}, \gamma N_{n+1}^{-\sigma}\right)$, $\mathcal{G}_{n+1}$ is a contraction in $\mathcal{B}_{n+1}:=\left\{h \in W^{\left(N_{n+1}\right)} \mid\|h\|_{s} \leq \rho_{n+1}:=2 K_{1} \varepsilon \gamma^{-1} N_{n+1}^{-\sigma-1}\right\}$.

Proof. For $N_{0}$ large enough (i.e. $\lambda$ large enough, see (8)) we have $N_{n+1}^{-\sigma}<N_{n}^{-\sigma} / 2, \forall n$, and so $\mathcal{N}\left(A_{n+1}, \gamma N_{n+1}^{-\sigma}\right) \subset \mathcal{N}\left(A_{n} ; \gamma N_{n}^{-\sigma} / 2\right)$. Then (19) holds and

$$
\left\|\mathcal{G}_{n+1}(h)\right\|_{s} \quad \stackrel{(21)}{\leq} \quad \frac{K}{\gamma} N_{n+1}^{\mu}\left(\left\|r_{n}\right\|_{s}+\left\|R_{n}(h)\right\|_{s}\right) .
$$

By (22) and the definition of $\beta$ in (16), if $\|h\|_{s} \leq \rho_{n+1}:=2 K_{1} \varepsilon \gamma^{-1} N_{n+1}^{-\sigma-1}$, then

$$
\left\|\mathcal{G}_{n+1}(h)\right\|_{s} \leq \frac{\varepsilon}{\gamma} K_{1} N_{n+1}^{-\sigma-1}+\frac{\varepsilon}{\gamma} K_{1} N_{n+1}^{\mu} \rho_{n+1}^{2} \leq \rho_{n+1}
$$

for $\varepsilon \gamma^{-1}$ small enough.
Next we use the formula

$$
D \mathcal{G}_{n+1}(h)[v]=\varepsilon \mathcal{L}_{N_{n+1}}^{-1}\left(\varepsilon, \omega, \widetilde{u}_{n}\right) P_{N_{n+1}}\left(\left(D f\left(\widetilde{u}_{n}+h\right)-D f\left(\widetilde{u}_{n}\right)\right) v\right)
$$

to obtain

$$
\begin{equation*}
\left\|D \mathcal{G}_{n+1}(h)[v]\right\|_{s} \leq K \frac{\varepsilon}{\gamma} N_{n+1}^{\mu}\|h\|_{s}\|v\|_{s} \leq K \frac{\varepsilon}{\gamma} N_{n+1}^{\mu} \rho_{n+1}\|v\|_{s} \leq \frac{\|v\|_{s}}{2} \tag{23}
\end{equation*}
$$

for $\varepsilon \gamma^{-1}$ small enough. Hence $\mathcal{G}_{n+1}$ is a contraction on $\mathcal{B}_{n+1}$.
Let $h_{n+1}$ be the fixed point of $\mathcal{G}_{n+1}$ defined for all $(\varepsilon, \omega) \in \mathcal{N}\left(A_{n+1}, \gamma N_{n+1}^{-\sigma}\right)$ and $u_{n+1}:=\widetilde{u}_{n}+h_{n+1}$.
Lemma 2.2 (Estimate of the derivatives) The map $h_{n+1} \in C^{1}\left(\mathcal{N}\left(A_{n+1}, \gamma N_{n+1}^{-\sigma}\right) ; W^{\left(N_{n+1}\right)}\right)$ and $\left\|D_{\varepsilon, \omega} h_{n+1}\right\|_{s} \leq K \gamma^{-1} N_{n+1}^{-1}$.
Proof. For all $(\varepsilon, \omega) \in \mathcal{N}\left(A_{n+1}, \gamma N_{n+1}^{-\sigma}\right), h_{n+1}(\varepsilon, \omega)$ is a solution of $U_{n+1}\left(\varepsilon, \omega, h_{n+1}(\varepsilon, \omega)\right)=0$ where

$$
U_{n+1}(\varepsilon, \omega, h):=L_{\omega}\left(\widetilde{u}_{n}+h\right)-\varepsilon P_{N_{n+1}} f\left(\widetilde{u}_{n}+h\right)
$$

Note that, with the notations of the proof of lemma 2.1, $U_{n+1}(\varepsilon, \omega, h)=\mathcal{L}_{N_{n+1}}\left(\varepsilon, \omega, \widetilde{u}_{n}\right)\left(h-\mathcal{G}_{n+1}(h)\right)$. Since $\left\|D \mathcal{G}_{n+1}\left(h_{n+1}\right)\right\|_{s} \leq 1 / 2$ (see (23)), $D_{h} U_{n+1}\left(\varepsilon, \omega, h_{n+1}\right)=\mathcal{L}_{N_{n+1}}\left(\varepsilon, \omega, u_{n+1}\right)$ is invertible and

$$
\begin{equation*}
\left\|\left(D_{h} U_{n+1}\left(\varepsilon, \omega, h_{n+1}\right)\right)^{-1}\right\|_{s}=\left\|\left(I-D \mathcal{G}_{n+1}\left(h_{n+1}\right)\right)^{-1} \mathcal{L}_{N_{n+1}}^{-1}\left(\varepsilon, \omega, \widetilde{u}_{n}\right)\right\|_{s} \leq \frac{K}{\gamma} N_{n+1}^{\mu} \tag{24}
\end{equation*}
$$

Hence, by the Implicit Function Theorem, $h_{n+1} \in C^{1}\left(\mathcal{N}\left(A_{n+1}, \gamma N_{n+1}^{-\sigma}\right), W^{\left(N_{n+1}\right)}\right)$ and

$$
\begin{equation*}
D_{\varepsilon, \omega} h_{n+1}=-\mathcal{L}_{N_{n+1}}^{-1}\left(\varepsilon, \omega, u_{n+1}\right)\left(D_{\varepsilon, \omega} U_{n+1}\right)\left(\varepsilon, \omega, h_{n+1}\right) . \tag{25}
\end{equation*}
$$

Now, using that $\widetilde{u}_{n}(\varepsilon, \omega)$ solves $\left(P_{N_{n}}\right)$ for $(\varepsilon, \omega) \in \mathcal{N}\left(A_{n+1}, \gamma N_{n+1}^{-\sigma}\right)$, we get

$$
\begin{equation*}
D_{\varepsilon} U_{n+1}\left(\varepsilon, \omega, h_{n+1}\right)=P_{N_{n}} f\left(\widetilde{u}_{n}\right)-P_{N_{n+1}} f\left(u_{n+1}\right)+\varepsilon P_{N_{n}} D f\left(\widetilde{u}_{n}\right) D_{\varepsilon} \widetilde{u}_{n}-\varepsilon P_{N_{n+1}} D f\left(u_{n+1}\right) D_{\varepsilon} \widetilde{u}_{n} \tag{26}
\end{equation*}
$$

and

$$
\begin{equation*}
D_{\omega} U_{n+1}\left(\varepsilon, \omega, h_{n+1}\right)=2 \omega\left(h_{n+1}\right)_{t t}+\varepsilon P_{N_{n}} D f\left(\widetilde{u}_{n}\right) D_{\omega} \widetilde{u}_{n}-\varepsilon P_{N_{n+1}} D f\left(u_{n+1}\right) D_{\omega} \widetilde{u}_{n} . \tag{27}
\end{equation*}
$$

We deduce from (24)-(27) the estimate $\left\|D_{\varepsilon, \omega} h_{n+1}\right\|_{s} \leq K \gamma^{-1} N_{n+1}^{-1}$ using (11), (17), (18), $\left\|h_{n+1}\right\|_{s} \leq \rho_{n+1}$ in lemma 2.1, and the definition of $\beta$ in (16); we omit the details.

We now define, by means of a cut-off function, a $C^{1}$-extension of $h_{n+1} \in C^{1}\left(A_{n+1}, W^{\left(N_{n}\right)}\right)$ onto the whole $A$. Note that $h_{n+1}$ is yet defined on a neighborhood of $A_{n+1}$ of width $\gamma N_{n+1}^{-\sigma}$ and $\left\|h_{n+1}\right\|_{s}=$ $O\left(\varepsilon \gamma^{-1} N_{n+1}^{-\sigma-1}\right)$.

Lemma 2.3 (Whitney extension) There exists an extension $\widetilde{h}_{n+1} \in C^{1}\left(A, W^{\left(N_{n+1}\right)}\right)$ of $h_{n+1}$ satisfying, for $\varepsilon \gamma^{-1}$ small, $\left\|\widetilde{h}_{n+1}\right\|_{s} \leq K \varepsilon \gamma^{-1} N_{n+1}^{-\sigma-1},\left\|D_{\varepsilon, \omega} \widetilde{h}_{n+1}\right\|_{s} \leq K \gamma^{-1} N_{n+1}^{-1}$.

Proof. Let

$$
\widetilde{h}_{n+1}(\varepsilon, \omega):=\left\{\begin{array}{lll}
\psi_{n+1}(\varepsilon, \omega) h_{n+1}(\varepsilon, \omega) & \text { if } & (\varepsilon, \omega) \in \mathcal{N}\left(A_{n+1}, \gamma N_{n+1}^{-\sigma}\right)  \tag{28}\\
0 & \text { if } & (\varepsilon, \omega) \notin \mathcal{N}\left(A_{n+1}, \gamma N_{n+1}^{-\sigma}\right)
\end{array}\right.
$$

where $\psi_{n+1}$ is a $C^{\infty}$ cut-off function satisfying $0 \leq \psi_{n+1} \leq 1, \psi_{n+1} \equiv 1$ on $\mathcal{N}\left(A_{n+1}, \gamma N_{n+1}^{-\sigma} / 2\right), \psi_{n+1} \equiv 0$ outside $\mathcal{N}\left(A_{n+1}, \gamma N_{n+1}^{-\sigma}\right)$, and $\left|D_{\varepsilon, \omega} \psi_{n+1}\right| \leq \gamma^{-1} N_{n+1}^{\sigma} C$ (a cut-off function $\psi_{n+1}$ can be constructed like in lemma 3.3 of [3]).

Then $\left\|\widetilde{h}_{n+1}\right\|_{s} \leq\left\|h_{n+1}\right\|_{s} \leq K \varepsilon \gamma^{-1} N_{n+1}^{-\sigma-1}$ by lemma 2.1, and, for $\varepsilon \gamma^{-1} \leq 1$,

$$
\left\|D_{\varepsilon, \omega} \widetilde{h}_{n+1}\right\|_{s} \leq\left|D_{\varepsilon, \omega} \psi_{n+1}\right|\left\|h_{n+1}\right\|_{s}+\left\|D_{\varepsilon, \omega} h_{n+1}\right\|_{s} \leq \bar{K} \gamma^{-1} N_{n+1}^{-1}
$$

thanks to the estimate $\left\|D_{\varepsilon, \omega} h_{n+1}\right\|_{s} \leq K \gamma^{-1} N_{n+1}^{-1}$ proved in lemma 2.2.
Finally we define the $C^{1}$-Whitney extension $\widetilde{u}_{n+1} \in C^{1}\left(A, W^{\left(N_{n+1}\right)}\right)$ of $u_{n+1}$ as

$$
\widetilde{u}_{n+1}:=\widetilde{u}_{n}+\widetilde{h}_{n+1} .
$$

To complete the induction proof of Theorem 2.1 we have to prove that $\widetilde{u}_{n+1}$ and $A_{n+1}$ satisfy all the properties stated in the Theorem plus (14)-(15) at the step $n+1$.

By lemma 2.3 property (14) is verified at the step $n+1$. Then we get also $\left\|\widetilde{u}_{n+1}\right\|_{s} \leq K \varepsilon \gamma^{-1}$ and $\left\|D_{\varepsilon, \omega} \widetilde{u}_{n+1}\right\|_{s} \leq K \gamma^{-1}$.

Now we prove that also (15) holds at the step $n+1$.
Lemma 2.4 For $\varepsilon \leq \varepsilon_{0}\left(\gamma, N_{0}, \beta\right)$ small enough,

$$
B_{n+1} \leq\left(1+N_{n+1}^{\mu}\right) B_{n} \quad \text { and } \quad B_{n+1}^{\prime} \leq B_{n}^{\prime}+K \gamma^{-1} N_{n+1}^{\mu}\left(N_{n+1}^{\sigma} B_{n}+\varepsilon B_{n}^{\prime}\right)
$$

Proof. We have $B_{n+1} \leq B_{n}+\left\|\widetilde{h}_{n+1}\right\|_{s+\beta}$ where, by (28), \| $\left\|\widetilde{h}_{n+1}\right\|_{s+\beta} \leq\left\|h_{n+1}\right\|_{s+\beta}$ and, for all $(\varepsilon, \omega) \in \mathcal{N}\left(A_{n+1}, \gamma N_{n+1}^{-\sigma}\right), h_{n+1}=-\mathcal{L}_{N_{n+1}}^{-1}\left(\varepsilon, \omega, \widetilde{u}_{n}\right)\left(r_{n}+R_{n}\left(h_{n+1}\right)\right)$. Then

$$
\begin{equation*}
\left\|h_{n+1}\right\|_{s+\beta} \stackrel{(P 5)}{\leq} K \frac{N_{n+1}^{\mu}}{\gamma}\left(\left\|r_{n}\right\|_{s+\beta}+\left\|R_{n}\left(h_{n+1}\right)\right\|_{s+\beta}+\left\|\widetilde{u}_{n}\right\|_{s+\beta}\left(\left\|r_{n}\right\|_{s}+\left\|R_{n}\left(h_{n+1}\right)\right\|_{s}\right)\right) \tag{29}
\end{equation*}
$$

By (19) and the tame estimate (P2),

$$
\begin{equation*}
\left\|r_{n}\right\|_{s+\beta} \leq \varepsilon\left\|f\left(\widetilde{u}_{n}\right)\right\|_{s+\beta} \stackrel{(P 2)}{\leq} \varepsilon K\left(1+\left\|\widetilde{u}_{n}\right\|_{s+\beta}\right)=\varepsilon K B_{n} \tag{30}
\end{equation*}
$$

By the Taylor tame estimate (P3), and since $\left\|h_{n+1}\right\|_{s} \leq \rho_{n+1}=2 K_{1} \varepsilon \gamma^{-1} N_{n+1}^{-\sigma-1}$ (lemma 2.1)

$$
\begin{equation*}
\left\|R_{n}\left(h_{n+1}\right)\right\|_{s+\beta} \stackrel{(P 3)}{\leq} \varepsilon K\left(\left\|\widetilde{u}_{n}\right\|_{s+\beta}\left\|h_{n+1}\right\|_{s}^{2}+\left\|h_{n+1}\right\|_{s}\left\|h_{n+1}\right\|_{s+\beta}\right) \leq \varepsilon K\left(B_{n} \rho_{n+1}^{2}+\rho_{n+1}\left\|h_{n+1}\right\|_{s+\beta}\right) \tag{31}
\end{equation*}
$$

Inserting in (29) estimates (30)-(31), $\left\|r_{n}\right\|_{s},\left\|R_{n}\left(h_{n+1}\right)\right\|_{s} \leq K \varepsilon$, we get

$$
\left\|h_{n+1}\right\|_{s+\beta} \leq \bar{K} \frac{\varepsilon}{\gamma} N_{n+1}^{\mu} B_{n}+\left(\bar{K} \frac{\varepsilon}{\gamma} N_{n+1}^{\mu-\sigma-1}\right)\left\|h_{n+1}\right\|_{s+\beta} \leq \bar{K} \frac{\varepsilon}{\gamma} N_{n+1}^{\mu} B_{n}+\frac{1}{2}\left\|h_{n+1}\right\|_{s+\beta}
$$

for $\bar{K} \varepsilon \gamma^{-1}<1 / 2$. Hence $\left\|h_{n+1}\right\|_{s+\beta} \leq N_{n+1}^{\mu} B_{n}$. The second inequality follows similarly by (25)-(27) and using $\left\|D_{\varepsilon, \omega} \widetilde{h}_{n+1}\right\|_{s+\beta} \leq C \gamma^{-1} N_{n+1}^{\sigma}\left\|h_{n+1}\right\|_{s+\beta}+\left\|D_{\varepsilon, \omega} h_{n+1}\right\|_{s+\beta}$.

Lemma 2.5 (Measure estimate) $A_{\infty}:=\bigcap_{n=0}^{\infty}$ satisfies $\left|A_{\infty}\right| \geq \varepsilon_{0}\left(\bar{\omega}_{2}-\bar{\omega}_{1}\right)-C \varepsilon_{0} \gamma$.
Proof. Recalling (20), the complementary $A_{\infty}^{c}=\bigcup_{n=0}^{\infty} G_{N_{n+1}}^{c}\left(\widetilde{u}_{n}\right)$ and, by (P5), $G_{N_{n+1}}^{c}\left(\widetilde{u}_{n}\right)=G^{c}$ for all $n$ such that $N_{n+1} \leq c\left(\gamma \varepsilon_{0}^{-1}\right)^{2 / 3}$, that is all $n=0,1 \ldots, n^{*}$ for some $n^{*}$, so that $A_{n^{*}+1}=A_{n^{*}}=\ldots=A_{0}$. By standard arguments, $\left|G^{c}\right| \leq C \varepsilon_{0} \gamma$. Since $\left\|\widetilde{u}_{n}-\widetilde{u}_{n-1}\right\| \leq N_{n}^{-\sigma}, \forall n$, by (13),

$$
\begin{aligned}
\left|\bigcup_{n=0}^{\infty} G_{N_{n+1}}^{c}\left(\widetilde{u}_{n}\right)\right| & =\sum_{n>n^{*}}^{\infty}\left|G_{N_{n+1}}^{c}\left(\widetilde{u}_{n}\right) \backslash G_{N_{n}}^{c}\left(\widetilde{u}_{n-1}\right)\right|+\left|G^{c}\right| \\
& \leq \sum_{n>n^{*}} \varepsilon_{0} \gamma N_{n}^{-1}+\left|G^{c}\right| \leq C^{\prime} \varepsilon_{0} \gamma
\end{aligned}
$$

and we get the thesis.
Finally, by (14), we deduce that, $\forall(\varepsilon, \omega) \in A_{\infty}$, the series

$$
\widetilde{u}:=\widetilde{u}_{0}+\sum_{n=1}^{\infty}\left(\widetilde{u}_{n}-\widetilde{u}_{n-1}\right)
$$

converges in $H^{s}$ (uniformly in $(\varepsilon, \omega)$ ) to a solution of equation (3) and $\|\widetilde{u}\|_{s} \leq K \varepsilon \gamma^{-1},\left\|D_{\varepsilon, \omega} \widetilde{u}\right\|_{s} \leq K \gamma^{-1}$. The proof of Theorem 2.1 is complete.

## 3 The linearized problem: proof of (P5)

### 3.1 Preliminaries

For $A \subset \mathbf{Z}^{d+1}$, $A$ finite and symmetric (i.e. $-A=A$ ), we define the finite dimensional subspace of $H^{s}$

$$
H_{A}:=\operatorname{Span}_{k \in A}\left\{e_{k}\right\}=\left\{\sum_{k \in A} h_{k} e_{k}: h_{k} \in \mathbf{C}, h_{k}^{*}=h_{-k}\right\} \quad \text { where } \quad e_{k}(t, x):=e^{\mathrm{i}(l t+j \cdot x)}
$$

and $k:=(l, j) \in \mathbf{Z} \times \mathbf{Z}^{d}$. We shall denote by $P_{A}$ the $L^{2}$-orthogonal projector on $H_{A}$, defined by

$$
P_{A} h:=\sum_{k \in A} h_{k} e_{k}, \quad \forall h=\sum_{k \in \mathbf{Z}^{d+1}} h_{k} e_{k} \in H^{s}
$$

Note that $P_{A}$ is also the $H^{s}$-orthogonal projector onto $H_{A}$. With these notations, the linear operator $\mathcal{L}_{N}$ in (10) is defined on $H_{\Omega_{N}}$, where

$$
\Omega_{N}:=\left\{k:=(l, j) \in \mathbf{Z} \times \mathbf{Z}^{d}|\quad|(l, j) \mid \leq N\right\}
$$

and we write ${ }^{5}$

$$
\mathcal{L}_{N}=D+\varepsilon T \quad \text { with } \quad D h:=\omega^{2} h_{t t}-\Delta h+m h, \quad T h:=P_{\Omega_{N}}(a h), \quad a(t, x):=-\left(\partial_{u} F\right)(t, x, u(t, x)) .
$$

Note that $T$, as $u$, depends on the parameters $\varepsilon$ and $\omega$. In the $L^{2}$ - orthonormal basis $\left(e_{k}\right)_{k \in \Omega_{N}}$ of $H_{\Omega_{N}}$, $D$ is represented by a diagonal matrix with eigenvalues

$$
d_{k}=d_{(l, j)}:=-\omega^{2} l^{2}+\lambda_{j}, \quad \lambda_{j}:=\left|j^{2}\right|+m
$$

whereas $T$ is represented by the self-adjoint Toepliz matrix $\left(a_{k-k^{\prime}}\right)_{k, k^{\prime} \in \Omega_{N}}$, the $a_{k}$ being the Fourier coefficients of the function $a(t, x)$.

Definition 3.1 Given $U \in \mathcal{L}\left(H_{\Omega_{N}}, H_{\Omega_{N}}\right)$ and $A, B \subset \Omega_{N}$, we define the linear operator

$$
U_{B}^{A}: H_{A} \rightarrow H_{B}, \quad U_{B}^{A}:=\left.P_{B} U\right|_{H_{A}}
$$

and its operatorial norm $\left\|U_{B}^{A}\right\|_{s}:=\sup \left\{\left\|U_{B}^{A} h\right\|_{s}, h \in H_{A},\|h\|_{s}=1\right\}$. We set

$$
L_{B}^{A}:=\left.P_{B} \mathcal{L}_{N}\right|_{H_{A}}=\left.P_{B}(D+\varepsilon T)\right|_{H_{A}}=D_{B}^{A}+\varepsilon T_{B}^{A}
$$

Note that $L_{A}^{B}$ is the $L^{2}$-adjoint of $L_{B}^{A}$ and, in particular, $L_{A}^{A}$ is $L^{2}$-selfadjoint.
Lemma 3.1 Let $U \in \mathcal{L}\left(H_{A}, H_{A}\right)$ be $L^{2}$-selfadjoint.
i) Its eigenvalues $\lambda_{k}$ satisfy $\min _{k}\left|\lambda_{k}\right| \geq \bar{\lambda}>0$ if and only if $\left\|U^{-1}\right\|_{0} \leq \bar{\lambda}^{-1}$.
ii) If true then, $\forall s^{\prime} \geq 0, \quad\left\|U^{-1}\right\|_{s^{\prime}} \leq \bar{\lambda}^{-1}\left(\frac{M_{A}}{m_{A}}\right)^{s^{\prime}} \quad$ where $\quad M_{A}:=\max _{k \in A}|k|, m_{A}:=\min _{k \in A}|k|$.

Proof. $i$ ) is standard. For $i i$ ), by the smoothing estimates $(P 4)$ and lemma 3.1 we get

$$
\left\|U^{-1} h\right\|_{s^{\prime}} \leq M_{A}^{s^{\prime}}\left\|U^{-1} h\right\|_{0} \leq \frac{M_{A}^{s^{\prime}}}{\bar{\lambda}}\|h\|_{0} \leq \frac{M_{A}^{s^{\prime}}}{m_{A}^{s^{\prime}}} \frac{1}{\bar{\lambda}}\|h\|_{s^{\prime}}
$$

A bounded subset $A$ of $\mathbf{Z}^{d+1}$ is said dyadic if $M_{A} \leq 2 m_{A}$ where $M_{A}:=\max _{k \in A}|k|, m_{A}:=\min _{k \in A}|k|$. By lemma 3.1, the inverse of the minimum modulus of the eigenvalues of a self-adjoint operator $U$, acting on a subspace $H_{A}$ with dyadic $A$, provides also a bound of the Sobolev operatorial norm $\left\|U^{-1}\right\|_{s^{\prime}}$.

In the next lemma we estimate the variations of the eigenvalues of $L_{A}^{A}(\varepsilon, \omega)$ with respect to $\omega$.
Lemma 3.2 Let $A \subset \Omega_{N}, 0<\bar{\omega}_{1}<\bar{\omega}_{2}$ and $I$ be any compact interval in $[-\gamma, \gamma]$, of length $|I|$. Suppose $\|u\|_{C^{1}(s)} \leq C \gamma^{-1}$. There exists $\varepsilon_{0}(\gamma)>0$ such that, for all $0<\varepsilon \leq \varepsilon_{0}(\gamma)$,

$$
\begin{equation*}
\mid\left\{\omega \in\left[\bar{\omega}_{1}, \bar{\omega}_{2}\right] \text { s.t. at least one eigenvalue of } L_{A}^{A}(\varepsilon, \omega) \text { belongs to } I\right\} \left\lvert\, \leq \frac{|A||I|}{\bar{\omega}_{1}}\right. \tag{32}
\end{equation*}
$$

where $|A|$ denotes the number of elements of $A$.
Proof. Since the map $\omega \mapsto L_{A}^{A}(\varepsilon, \omega)$ is $C^{1}$ and each $L_{A}^{A}(\varepsilon, \omega)$ is selfadjoint, its eigenvalues can be listed as $C^{1}$ functions $\lambda_{k}(\varepsilon, \omega), 1 \leq k \leq|A|$, of $\omega$. In what follows, $0<\varepsilon \leq \varepsilon_{0}(\gamma)$ is fixed and we omit the dependency with respect to $\varepsilon$.

Denoting $E_{\omega, k}$ the eigenspace of $L_{A}^{A}(\omega)$ associated to $\lambda_{k}(\omega)$,

$$
\begin{align*}
\left(\partial_{\omega} \lambda_{k}\right)(\omega) \leq \max _{h \in E_{\omega, k},\|h\|_{0}=1}\left(\left(\partial_{\omega} L_{A}^{A}\right)(\omega) h, h\right)_{0} & =\max _{h \in E_{\omega, k},\|h\|_{0}=1}\left(2 \omega\left(\partial_{t t} h, h\right)_{0}+O\left(\varepsilon \gamma^{-1}\right)\right) \\
& \leq \max _{h \in E_{\omega, k},\|h\|_{0}=1}\left(-2 \omega\left\|\partial_{t} h\right\|_{0}^{2}+O\left(\varepsilon \gamma^{-1}\right)\right) \tag{33}
\end{align*}
$$

[^4]using that $\left\|\partial_{\omega} a\right\|_{s}=\left\|\left(\partial_{u}^{2} F\right)(t, x, u) \partial_{\omega} u\right\|_{s} \leq C \gamma^{-1}$ and integrating by parts.
Let $P_{1}$ denote the $L^{2}$-orthogonal projector on $H_{1, A}:=\left\{h \in H_{A} \mid \partial_{t} h=0\right\}$.
Claim: Let $\lambda_{k}(\omega) \in[-\gamma, \gamma]$. For $h \in E_{\omega, k}$, decompose $h=h_{1}+h_{2}$ with $h_{1}:=P_{1} h, h_{2}:=\left(I-P_{1}\right) h \in H_{1, A}^{\perp}$. Then, for $\varepsilon_{0}$ small enough, $\left\|h_{2}\right\|_{0}^{2} \geq 3\|h\|_{0}^{2} / 4$.
Indeed, if $h \in E_{\omega, k}, \omega^{2} h_{t t}-\Delta h+m h+\varepsilon P_{A}(a h)=\lambda_{k}(\omega) h$ and so $h_{1}$ satisfies
\[

$$
\begin{equation*}
-\Delta h_{1}+m h_{1}-\lambda_{k}(\omega) h_{1}=-\varepsilon P_{1} P_{A}(a h) . \tag{34}
\end{equation*}
$$

\]

Furthermore, by (2), $\left||j|^{2}+m\right| \geq 3 \gamma, \forall j \in \mathbf{Z}^{d}$, and if $\lambda_{k}(\omega) \in[-\gamma, \gamma]$ then $|\| j|^{2}+m-\lambda_{k}(\omega) \mid \geq 2 \gamma$. Hence $-\Delta+m-\lambda_{k}(\varepsilon, \omega)$ is a $L^{2}$-selfadjoint operator of $H_{A}$ whose eigenvalues are of modulus $\geq 2 \gamma$. From (34),

$$
\left\|h_{1}\right\|_{0} \leq \frac{\varepsilon}{2 \gamma}\left\|P_{1} P_{A}(a h)\right\|_{0} \leq \frac{C \varepsilon}{\gamma}\|a\|_{s}\|h\|_{0} \leq \frac{\|h\|_{0}}{2}
$$

for $\varepsilon \gamma^{-1}$ small enough, and the claim follows.
Finally, since $\left\|\partial_{t} h\right\|_{0}^{2}=\left\|\partial_{t} h_{2}\right\|_{0}^{2} \geq\left\|h_{2}\right\|_{0}^{2}$, we deduce, from (33) and the previous claim,

$$
\left(\partial_{\omega} \lambda_{k}\right)(\omega) \leq \max _{h \in E_{\omega, k},\|h\|_{0}=1}\left(-2 \omega\left\|h_{2}\right\|_{0}^{2}+O\left(\varepsilon \gamma^{-1}\right)\right) \leq-\frac{3 \omega}{2}+O\left(\varepsilon \gamma^{-1}\right) \leq-\omega \leq-\bar{\omega}_{1}
$$

for $0<\varepsilon \leq \varepsilon_{0}(\gamma)$ small enough. Therefore, $\left|\lambda_{k}^{-1}(I) \cap\left(\bar{\omega}_{1}, \bar{\omega}_{2}\right)\right| \leq|I| / \bar{\omega}_{1}$ and, summing over all the $|A|$ eigenvalues, we deduce (32).

### 3.2 Regular and Singular sites

Definition 3.2 Fixed $\rho>0$ we define the regular sites $R$ and the singular sites $S$, as

$$
R:=\left\{k \in \Omega_{N}| | d_{k} \mid \geq \rho\right\} \quad \text { and } \quad S:=\Omega_{N} \backslash R:=\left\{k \in \Omega_{N}| | d_{k} \mid<\rho\right\} .
$$

We introduce the following assumption: the singular sites $S$ can be partitioned in disjoint clusters $\Omega_{\alpha}$,

$$
\begin{equation*}
S=\bigcup_{\alpha \in \mathcal{A}} \Omega_{\alpha} \tag{35}
\end{equation*}
$$

satisfying:

- (H1) (dyadic) $M_{\alpha} \leq 2 m_{\alpha}, \forall \alpha \in \mathcal{A}$, where $M_{\alpha}:=\max _{k \in \Omega_{\alpha}}|k|, m_{\alpha}:=\min _{k \in \Omega_{\alpha}}|k|$.
- (H2) (separation) $\exists \delta:=\delta(d) \in(0,1)$ such that $\mathrm{d}\left(\Omega_{\alpha}, \Omega_{\sigma}\right):=\min _{k \in \Omega_{\alpha}, k^{\prime} \in \Omega_{\sigma}}\left|k-k^{\prime}\right| \geq\left(M_{\alpha}+M_{\sigma}\right)^{\delta} /$ $2, \forall \alpha \neq \sigma$.

For the wave equation (3), we shall verify both properties (H1)-(H2) in lemma 3.6 and section 4.
Remark 3.1 Condition (H1) is weaker than the corresponding one in [7] (Lemma 7 of Appendix 2). It could be considerably weakened, but not completely eliminated, see remark 3.2.

For $K \leq N$ let

$$
\begin{equation*}
\Omega_{K}:=\left\{k \in \mathbf{Z}^{d+1}| | k \mid \leq K\right\} \quad \text { and } \quad L_{K}:=L_{\Omega_{K}}^{\Omega_{K}} \tag{36}
\end{equation*}
$$

Proposition 3.1 Assume that there is a partition of $S$ satisfying (H1)-(H2). Given $\tau>0$, there exist $\mu:=\mu(d, \tau), s:=s(\delta, d, \tau)$ (independent of $N$ ) and, $\forall \gamma \in(0,1), \forall s^{\prime} \geq s$, there is $\varepsilon_{0}:=\varepsilon_{0}\left(\gamma, s^{\prime}\right)>0$ such that : if $<\varepsilon \leq \varepsilon_{0}$ and $\forall K \leq N$, all the eigenvalues of $L_{K}$ have modulus $\geq \gamma K^{-\tau}$ then

$$
\begin{equation*}
\left\|\mathcal{L}_{N}^{-1} h\right\|_{s^{\prime}} \leq \frac{K\left(s^{\prime}\right)}{\gamma} N^{\mu}\left(\|h\|_{s^{\prime}}+\|u\|_{s^{\prime}}\|h\|_{s}\right), \quad \forall h \in H_{\Omega_{N}} \tag{37}
\end{equation*}
$$

Remark 3.2 Without exploiting any separation property of the clusters we would only have an estimate like $\forall N,\left\|\mathcal{L}_{N}^{-1} h\right\|_{s^{\prime}} \leq N^{s^{\prime}}\|h\|_{s^{\prime}}, \forall h \in H_{\Omega_{N}}$. This is even weaker than a tame like estimate like $\left\|\mathcal{L}_{N}^{-1} h\right\|_{s^{\prime}} \leq$ $K\|h\|_{2 s^{\prime}}$. In [19] Lojasiewicz and Zehnder have shown with a counter-example that a so weak condition is not enough for the convergence of the Nash-Moser iteration scheme, whereas a tame estimate like $\left\|\mathcal{L}_{N}^{-1} h\right\|_{s^{\prime}} \leq K\|h\|_{\lambda s^{\prime}}$ with $\lambda<2$ would be enough. This fact confirms that some more information, in addition to lower bounds for the moduli of the eigenvalues, is required.

We shall deduce property (P5) from Proposition 3.1 in the next section. The proof of Proposition 3.1 is provided in sections 3.4, 3.5, 3.6.

### 3.3 Proof of (P5)

Let

$$
B_{K}(u):=\left\{(\varepsilon, \omega) \in\left[0, \varepsilon_{0}\right] \times\left[\bar{\omega}_{1}, \bar{\omega}_{2}\right] \left\lvert\,\left\|L_{K}^{-1}\right\|_{0} \leq \frac{K^{\tau}}{2 \gamma}\right.\right\} \quad \text { and } \quad G_{N}(u):=\bigcap_{K \leq N} B_{K}(u) \cap G
$$

By lemma 3.1, $(\varepsilon, \omega) \in B_{K}(u)$ if and only if all the eigenvalues $\lambda_{k}, k=1, \ldots, K$, of $L_{K}$ have modulus greater or equal to $2 \gamma / K^{\tau}$. Below, to emphasize the dependence of $L_{K}$ with respect to parameters, we shall use the notations $L_{K}(\varepsilon, \omega)$.
Lemma 3.3 Let $\tau \geq 3 / 2$. There is $c>0$ such that $\forall K \leq N_{\varepsilon_{0}}:=c\left(\gamma \varepsilon_{0}^{-1}\right)^{2 / 3}$ we have $G \subset B_{K}(u)$. As a consequence, if $N \leq N_{\varepsilon_{0}}$ then $G_{N}(u)=G$.
Proof. The eigenvalues of $L_{K}$ have the form $-\omega^{2} l^{2}+|j|^{2}+m+O(\varepsilon)$ with $|(l, j)| \leq K$. Hence, $\forall(\varepsilon, \omega) \in G$, as a consequence of (2), if $\varepsilon_{0} \gamma^{-1} K^{3 / 2}$ is small enough, then all the eigenvalues of $L_{K}(\varepsilon, \omega)$ have modulus $\geq 2 \gamma / K^{3 / 2}$.

On the other hand, for $N>N_{\varepsilon_{0}}:=c\left(\gamma \varepsilon_{0}^{-1}\right)^{2 / 3}$, we have to excise additional "resonant parameters".
Lemma 3.4 Let $N>N_{\varepsilon_{0}}$. If $(\varepsilon, \omega) \in \mathcal{N}\left(G_{N}(u), \gamma N^{-3-\tau}\right)$ then, $\forall K \leq N$, all the eigenvalues of $L_{K}(\varepsilon, \omega)$ have modulus $\geq \gamma K^{-\tau}$.

Proof. For $\omega, \omega^{\prime} \in\left[\bar{\omega}_{1}, \bar{\omega}_{2}\right]$ and $\varepsilon, \varepsilon^{\prime} \in\left[0, \varepsilon_{0}\right]$, we have, $\forall h \in H_{\Omega_{K}}$,

$$
\begin{aligned}
\left\|\left(L_{K}(\varepsilon, \omega)-L_{K}\left(\varepsilon^{\prime}, \omega^{\prime}\right)\right) h\right\|_{0} & \leq\left|\omega^{2}-\omega^{\prime 2}\right|\left\|h_{t t}\right\|_{0}+\varepsilon^{\prime}\left\|a(\varepsilon, \omega)-a\left(\varepsilon^{\prime}, \omega^{\prime}\right)\right\|_{s}\|h\|_{0}+\left|\varepsilon-\varepsilon^{\prime}\right|\|a(\varepsilon, \omega)\|_{s}\|h\|_{0} \\
& \leq C K^{2}\left(\left|\omega-\omega^{\prime}\right|+\left|\varepsilon-\varepsilon^{\prime}\right|\right)\|h\|_{0}
\end{aligned}
$$

If $(\varepsilon, \omega) \in \mathcal{N}\left(G_{N}(u), \gamma N^{-3-\tau}\right)$, there is $\left(\varepsilon^{\prime}, \omega^{\prime}\right) \in G_{N}(u)$ such that $\left|\omega-\omega^{\prime}\right|+\left|\varepsilon-\varepsilon^{\prime}\right| \leq \gamma N^{-3-\tau}$ and $\forall K \leq N, \forall h \in H_{\Omega_{K}}$,

$$
\left\|L_{K}(\varepsilon, \omega) h\right\|_{0} \geq\left\|L_{K}\left(\varepsilon^{\prime}, \omega^{\prime}\right) h\right\|_{0}-C \frac{\gamma K^{2}}{N^{3+\tau}}\|h\|_{0} \geq\left(\frac{2 \gamma}{K^{\tau}}-C \frac{\gamma K^{2}}{N^{3+\tau}}\right)\|h\|_{0} \geq \frac{\gamma}{K^{\tau}}\|h\|_{0}
$$

by the definition of $B_{K}(u)$. The result follows.
Lemma 3.5 Let $\sigma, \tau \geq d+3$. Then the measure estimate (13) holds.
Proof. For $N^{\prime} \geq N$,

$$
G_{N^{\prime}}^{c}\left(u_{2}\right) \backslash G_{N}^{c}\left(u_{1}\right)=G_{N^{\prime}}^{c}\left(u_{2}\right) \cap G_{N}\left(u_{1}\right) \subset\left[\cup_{K \leq N}\left(B_{K}^{c}\left(u_{2}\right) \cap B_{K}\left(u_{1}\right) \cap G\right)\right] \cup\left[\cup_{K>N} B_{K}^{c}\left(u_{2}\right) \cap G\right] .
$$

By lemma 3.3, if $K \leq N_{\varepsilon_{0}}$ then $B_{K}^{c}\left(u_{2}\right) \cap G=\emptyset$. Hence, to show (13), it is enough to prove that, if $\left\|u_{1}-u_{2}\right\|_{s} \leq N^{-\sigma}$, then

$$
\sum_{K \leq N}\left|B_{K}^{c}\left(u_{2}\right) \cap B_{K}\left(u_{1}\right)\right|+\sum_{K>\max \left(N, N_{\varepsilon_{0}}\right)}\left|B_{K}^{c}\left(u_{2}\right)\right| \leq \bar{C} \frac{\gamma \varepsilon_{0}}{N}
$$

Since $\left\|L_{K}\left(u_{2}\right)-L_{K}\left(u_{1}\right)\right\|_{0}=O\left(\varepsilon\left\|u_{2}-u_{1}\right\|_{s}\right)=O\left(\varepsilon N^{-\sigma}\right)$, if one of the eigenvalues of $L_{K}\left(u_{2}\right)$ is in $\left[-2 \gamma K^{-\tau}, 2 \gamma K^{-\tau}\right]$ then, by the variational characterization of the eigenvalues of a selfadjoint operator, one of the eigenvalues of $L_{K}\left(u_{1}\right)$ is in $\left[-2 \gamma K^{-\tau}-C \varepsilon N^{-\sigma}, 2 \gamma K^{-\tau}+C \varepsilon N^{-\sigma}\right]$. As a result

$$
\begin{aligned}
B_{K}^{c}\left(u_{2}\right) \cap B_{K}\left(u_{1}\right) \subset & \left\{(\varepsilon, \omega) \mid \exists \text { at least one eigenvalue of } L_{K}\left(u_{1}\right)\right. \\
& \text { with modulus in } \left.\left[2 \gamma K^{-\tau}, 2 \gamma K^{-\tau}+C \varepsilon N^{-\sigma}\right]\right\} .
\end{aligned}
$$

Then, by lemma $3.2,\left|\left\{\omega \mid(\varepsilon, \omega) \in B_{K}^{c}\left(u_{2}\right) \cap B_{K}\left(u_{1}\right)\right\}\right| \leq C \varepsilon N^{-\sigma}\left|\Omega_{K}\right| / \bar{\omega}_{1}$ for each $\varepsilon \in\left(0, \varepsilon_{0}\right]$, whence

$$
\left|B_{K}^{c}\left(u_{2}\right) \cap B_{K}\left(u_{1}\right)\right| \leq C^{\prime} \varepsilon_{0}^{2}\left|\Omega_{K}\right| N^{-\sigma} \leq C^{\prime} \varepsilon_{0}^{2} K^{d+1} N^{-\sigma} .
$$

Moreover, still by lemma 3.2, $\left|B_{K}^{c}\left(u_{2}\right)\right| \leq C \varepsilon_{0}\left|\Omega_{K}\right| \gamma K^{-\tau} / \bar{\omega}_{1} \leq C^{\prime} \varepsilon_{0} \gamma K^{d+1-\tau}$. Hence, for $\sigma, \tau \geq d+3$,

$$
\begin{aligned}
& \sum_{K \leq N}\left|B_{K}^{c}\left(u_{2}\right) \cap B_{K}\left(u_{1}\right)\right|+\sum_{K>\max \left(N, N_{\varepsilon_{0}}\right)}\left|B_{K}^{c}\left(u_{2}\right)\right| \\
\leq & C \varepsilon_{0}^{2}\left(\sum_{K \leq N} K^{d+1}\right) N^{-\sigma}+C \varepsilon_{0} \gamma\left(\sum_{K>\max \left(N, N_{\varepsilon_{0}}\right)} K^{d+1-\tau}\right) \\
\leq & C \varepsilon_{0}^{2} N^{d+2-\sigma}+C^{\prime} \varepsilon_{0} \gamma\left(\max \left(N, N_{\varepsilon_{0}}\right)\right)^{d+2-\tau} \leq \bar{C} \gamma \varepsilon_{0} N^{-1},
\end{aligned}
$$

provided that $\varepsilon_{0}$ is small enough.

Lemma 3.6 (Separation in clusters) $\forall \gamma \in(0,1)$ let $\rho(\gamma):=\gamma^{3 r(d)+1} / 8$ as defined in lemma 4.4. There is $\varepsilon_{0}(\gamma)$ such that $\forall N \geq N_{\varepsilon_{0}}:=c\left(\gamma \varepsilon_{0}^{-1}\right)^{2 / 3}, \forall(\varepsilon, \omega) \in \mathcal{N}\left(G, \gamma N^{-\sigma}\right), \sigma \geq 3$, there exists a decomposition of the singular sites

$$
\begin{equation*}
S:=\left\{(l, j) \in \Omega_{N}| | \omega^{2} l^{2}-\lambda_{j} \mid<\rho:=\frac{\rho(\gamma)}{2}\right\}=\bigcup_{\alpha \in \mathcal{A}} \Omega_{\alpha} \tag{38}
\end{equation*}
$$

in clusters, satisfying both properties (H1)-(H2) for some $\delta:=\delta(d)$.
Proof. Let $\left(\varepsilon, \omega^{\prime}\right) \in G$. By lemma $4.4, \forall \gamma \in(0,1)$,

$$
S^{\prime}:=\left\{(l, j) \in \Omega_{N}| |\left(\omega^{\prime}\right)^{2} l^{2}-\lambda_{j} \mid<\rho(\gamma)\right\}=\bigcup_{\alpha \in \mathcal{A}} \Omega_{\alpha}^{\prime}
$$

satisfying (H1)-(H2) for some $\delta:=\delta(d)$. If $\varepsilon_{0}(\gamma)$ is small enough, then $\forall N \geq N_{\varepsilon_{0}}, \gamma N^{-\sigma} \leq \rho(\gamma) N^{-2} / 4 \bar{\omega}_{2}$, and, $\forall\left|\omega-\omega^{\prime}\right| \leq \gamma N^{-\sigma}$, we have $S \subset S^{\prime}$. Therefore $S=\bigcup_{\alpha \in \mathcal{A}} \Omega_{\alpha}$ where $\Omega_{\alpha}:=\Omega_{\alpha}^{\prime} \cap S$ satisfy (H1)-(H2).

Proof of Property (P5) completed. Let $(\varepsilon, \omega) \in \mathcal{N}\left(G_{N}, \gamma N^{-\sigma}\right)$ with $\sigma \geq 3+\tau$. Then, by lemma 3.6, there exists a partition of the singular sites $S$ satisfying assumptions (H1)-(H2). Furthermore, by lemma 3.4, $\forall K \leq N$, all the eigenvalues of $L_{K}$ have modulus $\geq \gamma K^{-\tau}$. Therefore, by Proposition 3.1, $\mathcal{L}_{N}(\varepsilon, \omega, u)$ is invertible and (11) holds (we can fix $\tau=d+3$ by lemma 3.5).

The measure estimate (13) is proved in lemma 3.5.
The remainder of this section is devoted to the proof of Proposition 3.1.

### 3.4 Reduction along the regular sites

We have to solve the linear system

$$
\begin{equation*}
\mathcal{L}_{N} h=b, \quad h, b \in H_{\Omega_{N}} . \tag{39}
\end{equation*}
$$

According to the splitting of the indexes $\Omega_{N}=R \cup S$ we decompose

$$
H_{\Omega_{N}}=H_{R} \oplus H_{S}
$$

in orthogonal subspaces. Writing the unique decomposition $h=h_{R}+h_{S}, b=b_{R}+b_{S}$, with $h_{R}, b_{R} \in H_{R}$, $h_{S}, b_{S} \in H_{S}$, problem (39) is equivalent to

$$
\left\{\begin{array}{l}
L_{R}^{R} h_{R}+L_{R}^{S} h_{S}=b_{R}  \tag{40}\\
L_{S}^{R} h_{R}+L_{S}^{S} h_{S}=b_{S}
\end{array}\right.
$$

We recall that $s>(d+1) / 2$ so that $H^{s}$ is a Banach algebra.
To prove the invertibility of $L_{R}^{R}=D_{R}+\varepsilon T_{R}$ where, for brevity, $D_{R}:=D_{R}^{R}, T_{R}:=T_{R}^{R}$, note that

$$
\begin{equation*}
\forall s \geq 0, \quad\left\|D_{R} h\right\|_{s} \geq\left(\min _{k \in R}\left|d_{k}\right|\right)\|h\|_{s} \geq \rho\|h\|_{s}, \quad \forall h \in H_{R} \tag{41}
\end{equation*}
$$

Lemma 3.7 There is $c>0$ such that if $\varepsilon\|a\|_{s} / \rho \leq c$ then $\left\|\left(L_{R}^{R}\right)^{-1}\right\|_{0},\left\|\left(L_{R}^{R}\right)^{-1}\right\|_{s} \leq 2 \rho^{-1}$.
Proof. We have, $\forall h \in H_{R}$,

$$
\left\|L_{R}^{R} h\right\|_{s} \geq\left\|D_{R} h\right\|_{s}-\varepsilon\left\|T_{R} h\right\|_{s} \geq\left(\rho-\varepsilon C\|a\|_{s}\right)\|h\|_{s}
$$

and the result follows. Similarly for $\left\|\left(L_{R}^{R}\right)^{-1}\right\|_{0}$.
Next we estimate the Sobolev norm $\left\|\left(L_{R}^{R}\right)^{-1} h\right\|_{s^{\prime}}$ for $s^{\prime} \geq s$. Since $T$ comes from the multiplication operator for the function $a \in H^{s}$, it is natural to expect the following interpolation type inequality.
Lemma 3.8 For $\varepsilon\left(\|a\|_{s}+1\right) / \rho \leq c\left(s^{\prime}\right)$ small enough, $\left\|\left(L_{R}^{R}\right)^{-1} h\right\|_{s^{\prime}} \leq 2 \rho^{-1}\left(\|a\|_{s^{\prime}}\|h\|_{s}+\|h\|_{s^{\prime}}\right)$.
Proof For $h \in H_{R}$ let $v:=\left(L_{R}^{R}\right)^{-1} h$. Then

$$
\|h\|_{s^{\prime}}=\left\|D_{R} v+\varepsilon T_{R} v\right\|_{s^{\prime}} \geq\left\|D_{R} v\right\|_{s^{\prime}}-\varepsilon\left\|T_{R} v\right\|_{s^{\prime}} \stackrel{(41)}{\geq} \rho\|v\|_{s^{\prime}}-\varepsilon\left\|T_{R} v\right\|_{s^{\prime}}
$$

Now $\left\|T_{R} h\right\|_{s^{\prime}} \leq C\left(s^{\prime}\right)\left(\|a\|_{s^{\prime}}\|h\|_{s}+\|a\|_{s}\|h\|_{s^{\prime}}\right), \forall h \in H_{R}$, (see e.g. [24]), hence

$$
\|h\|_{s^{\prime}} \geq\left(\rho-\varepsilon C\left(s^{\prime}\right)\|a\|_{s}\right)\|v\|_{s^{\prime}}-\varepsilon C\left(s^{\prime}\right)\|a\|_{s^{\prime}}\|v\|_{s}
$$

Therefore, if $\varepsilon C\left(s^{\prime}\right)\left(\|a\|_{s}+1\right) \rho^{-1} \leq 1 / 2$,

$$
\|v\|_{s^{\prime}} \leq \frac{2}{\rho}\left(\|h\|_{s^{\prime}}+\varepsilon C\left(s^{\prime}\right)\|a\|_{s^{\prime}}\|v\|_{s}\right) \leq \frac{2}{\rho}\left(\|h\|_{s^{\prime}}+\|a\|_{s^{\prime}}\|h\|_{s}\right)
$$

by Lemma 3.7.
Lemma 3.9 $\forall A, B \subset \Omega_{N}, \forall s^{\prime} \geq s,\left\|T_{B}^{A}\right\|_{0} \leq \frac{C\left(s^{\prime}\right)\|a\|_{s^{\prime}}}{(1+\mathrm{d}(A, B))^{s^{\prime}-(d+1) / 2}}$.
Proof Using Hölder inequality and exchanging the order of integration

$$
\begin{aligned}
\left\|T_{B}^{A} h\right\|_{0}^{2} & =\sum_{k \in B}\left|\sum_{m \in A} a_{k-m} h_{m}\right|^{2} \leq \sum_{k \in B}\left(\sum_{m \in A}\left|a_{k-m}\right|^{2}\left|h_{m}\right|^{2}\left(1+|k-m|^{2 s^{\prime}}\right)\right)\left(\sum_{m \in A} \frac{1}{1+|k-m|^{2 s^{\prime}}}\right) \\
& \leq \sum_{m \in A}\left|h_{m}\right|^{2}\left(\sum_{k \in B}\left|a_{k-m}\right|^{2}\left(1+|k-m|^{2 s^{\prime}}\right)\right) \sum_{|m| \geq \Delta} \frac{1}{1+|m|^{2 s^{\prime}}} \leq \frac{C\left(s^{\prime}\right)\|h\|_{0}^{2}\|a\|_{s^{\prime}}^{2}}{(1+\Delta)^{2 s^{\prime}-(d+1)}}
\end{aligned}
$$

where $\Delta:=\mathrm{d}(A, B)$.
Lemma 3.10 Let $U \in \mathcal{L}\left(H_{R}, H_{R}\right)$ be such that, for some $\kappa>0, \forall A, B \subset R$,

$$
\left\|U_{B}^{A}\right\|_{0} \leq \frac{C_{1}}{(1+\mathrm{d}(A, B))^{\kappa}}
$$

Then, $\forall p \geq 1, \forall A, B \subset R$,

$$
\begin{equation*}
\left\|\left(U^{p}\right)_{B}^{A}\right\|_{0} \leq \frac{C_{1} p^{\kappa+1}\|U\|_{0}^{p-1}}{(1+\mathrm{d}(A, B))^{\kappa}} \tag{42}
\end{equation*}
$$

Proof We proceed by induction noting that, for $p=1$, (42) is just our assumption. Now assume that (42) holds up to order $p$. Define $\Delta:=\mathrm{d}(A, B)$ and, for some $\lambda \in(0,1)$ to be specified later, $B^{\prime}=\mathcal{N}(B, \lambda \Delta) \cap R, B^{\prime \prime}=R \backslash B^{\prime}$. We have

$$
\begin{aligned}
\left\|\left(U^{p+1}\right)_{B}^{A}\right\|_{0} & =\left\|\left(U^{p}\right)_{B}^{B^{\prime}} U_{B^{\prime}}^{A}+\left(U^{p}\right)_{B}^{B^{\prime \prime}} U_{B^{\prime \prime}}^{A}\right\|_{0} \leq\left\|U^{p}\right\|_{0}\left\|U_{B^{\prime}}^{A}\right\|_{0}+\left\|\left(U^{p}\right)_{B}^{B^{\prime \prime}}\right\|_{0}\|U\|_{0} \\
& \leq\|U\|_{0}^{p} \frac{C_{1}}{\left(1+\mathrm{d}\left(A, B^{\prime}\right)\right)^{\kappa}}+\frac{C_{1} p^{\kappa+1}\|U\|_{0}^{p-1}}{\left(1+\mathrm{d}\left(B^{\prime \prime}, B\right)\right)^{\kappa}}\|U\|_{0}
\end{aligned}
$$

by our hypothesis of induction. Now $\mathrm{d}\left(A, B^{\prime}\right) \geq(1-\lambda) \Delta$ and $\mathrm{d}\left(B^{\prime \prime}, B\right) \geq \lambda \Delta$. Hence

$$
\left\|\left(U^{p+1}\right)_{B}^{A}\right\|_{0} \leq C_{1}\|U\|_{0}^{p}\left(\frac{1}{(1+(1-\lambda) \Delta)^{\kappa}}+\frac{p^{\kappa+1}}{(1+\lambda \Delta)^{\kappa}}\right) \leq \frac{C_{1}\|U\|_{0}^{p}}{(1+\Delta)^{\kappa}}\left(\frac{1}{(1-\lambda)^{\kappa}}+\frac{p^{\kappa+1}}{\lambda^{\kappa}}\right) .
$$

Since $\min _{\lambda \in(0,1)}\left[(1-\lambda)^{-\kappa}+p^{\kappa+1} \lambda^{-\kappa}\right]=(p+1)^{\kappa+1}$ we have proved (42) up to order $p+1$.
Lemma $3.11 \forall A, B \subset R$ with $A \cap B=\emptyset$, and for $\varepsilon\|a\|_{s} \rho^{-1} \leq c$ small enough,

$$
\left\|\left[\left(L_{R}^{R}\right)^{-1}\right]_{B}^{A}\right\|_{0} \leq \frac{\varepsilon C\left(s^{\prime}\right)\|a\|_{s^{\prime}}}{\rho^{2} \mathrm{~d}(A, B)^{s^{\prime}-(d+1) / 2}}
$$

Proof By the Newmann series expansion

$$
\begin{equation*}
\left(L_{R}^{R}\right)^{-1}:=\left(D_{R}+\varepsilon T_{R}\right)^{-1}=\left(I+\sum_{p \geq 1}(-1)^{p}\left(\varepsilon D_{R}^{-1} T_{R}\right)^{p}\right) D_{R}^{-1} \tag{43}
\end{equation*}
$$

and since $A \cap B=\emptyset,\left[\left(L_{R}^{R}\right)^{-1}\right]_{B}^{A}=\sum_{p \geq 1}(-1)^{p} \varepsilon^{p}\left(U^{p}\right)_{B}^{A}\left[D_{R}^{-1}\right]_{A}^{A}$ where $U:=D_{R}^{-1} T_{R}$. Since $\left\|D_{R}^{-1}\right\|_{0} \leq \rho^{-1}$,

$$
\begin{equation*}
\left\|\left[\left(L_{R}^{R}\right)^{-1}\right]_{B}^{A}\right\|_{0} \leq \frac{1}{\rho} \sum_{p \geq 1} \varepsilon^{p}\left\|\left(U^{p}\right)_{B}^{A}\right\|_{0} . \tag{44}
\end{equation*}
$$

Now, by lemma 3.9, $\forall A, B \subset R$,

$$
\left\|U_{B}^{A}\right\|_{0}=\left\|\left[D_{R}^{-1}\right]_{B}^{B} T_{B}^{A}\right\|_{0} \leq \frac{1}{\rho}\left\|T_{B}^{A}\right\|_{0} \leq \frac{C\left(s^{\prime}\right)\|a\|_{s^{\prime}}}{\rho(1+\mathrm{d}(A, B))^{\kappa}}
$$

with $\kappa:=s^{\prime}-(d+1) / 2$, whence, by lemma 3.10,

$$
\left\|\left(U^{p}\right)_{B}^{A}\right\|_{0} \leq \frac{C\left(s^{\prime}\right)\|a\|_{s^{\prime}} p^{\kappa+1}\|U\|_{0}^{p-1}}{\rho(1+\mathrm{d}(A, B))^{\kappa}} \leq \frac{C\left(s^{\prime}\right)\|a\|_{s^{\prime}} p^{\kappa+1}\|a\|_{s}^{p-1}}{\rho^{p} \mathrm{~d}(A, B)^{\kappa}}
$$

using that $\|U\|_{0}=\left\|D_{R}^{-1} T_{R}\right\|_{0} \leq \rho^{-1}\|a\|_{s}$. By (44) we get

$$
\left\|\left[\left(L_{R}^{R}\right)^{-1}\right]_{B}^{A}\right\|_{0} \leq \frac{\varepsilon C\left(s^{\prime}\right)\|a\|_{s^{\prime}}}{\rho^{2} \mathrm{~d}(A, B)^{\kappa}}\left(\sum_{p \geq 1}\left[\frac{\varepsilon\|a\|_{s}}{\rho}\right]^{p-1} p^{\kappa+1}\right)
$$

proving the lemma.
Solving the first equation in (40) gives

$$
\begin{equation*}
h_{R}=\left(L_{R}^{R}\right)^{-1}\left(b_{R}-L_{R}^{S} h_{S}\right) \tag{45}
\end{equation*}
$$

and inserting into the second equation in (40) gives

$$
\begin{equation*}
\left[L_{S}^{S}-L_{S}^{R}\left(L_{R}^{R}\right)^{-1} L_{R}^{S}\right] h_{S}=b_{S}-L_{S}^{R}\left(L_{R}^{R}\right)^{-1} b_{R} \tag{46}
\end{equation*}
$$

Our main task is to invert the selfadjoint operator

$$
\begin{equation*}
U: H_{S} \rightarrow H_{S}, \quad U:=L_{S}^{S}-L_{S}^{R}\left(L_{R}^{R}\right)^{-1} L_{R}^{S} \tag{47}
\end{equation*}
$$

According to (35) we get the orthogonal decomposition $H_{S}=\oplus_{\alpha \in \mathcal{A}} H_{\Omega_{\alpha}}$ which induces a block decomposition for the operator

$$
U: \oplus_{\alpha \in \mathcal{A}} H_{\Omega_{\alpha}} \rightarrow \oplus_{\alpha \in \mathcal{A}} H_{\Omega_{\alpha}}, \quad U=\left(U_{\Omega_{\sigma}}^{\Omega_{\alpha}}\right)_{\alpha, \sigma \in \mathcal{A}}, \quad U_{\Omega_{\sigma}}^{\Omega_{\alpha}}=L_{\Omega_{\sigma}}^{\Omega_{\alpha}}-L_{\Omega_{\sigma}}^{R}\left(L_{R}^{R}\right)^{-1} L_{R}^{\Omega_{\alpha}}
$$

where $\alpha$ is the column index and $\sigma$ is the row index. We write

$$
U=\mathcal{D}+\mathcal{R} \quad \text { with } \quad \mathcal{D}:=\operatorname{diag}_{\alpha \in \mathcal{A}}\left(U_{\Omega_{\alpha}}^{\Omega_{\alpha}}\right) \quad \text { and } \quad \mathcal{R}:=\left(U_{\Omega_{\sigma}}^{\Omega_{\alpha}}\right)_{\alpha \neq \sigma}
$$

### 3.5 Estimates of the blocks

We always assume in the sequel that $\varepsilon\|a\|_{s} / \rho<c$. We first estimate the off-diagonal blocks $U_{\Omega_{\sigma}}^{\Omega_{\alpha}}, \alpha \neq \sigma$.
Lemma $3.12 \forall s^{\prime} \geq s, \forall \alpha \neq \sigma$

$$
\begin{equation*}
\left\|U_{\Omega_{\sigma}}^{\Omega_{\alpha}}\right\|_{0} \leq \frac{\varepsilon C\left(s^{\prime}\right)\|a\|_{s^{\prime}}}{\mathrm{d}\left(\Omega_{\alpha}, \Omega_{\sigma}\right)^{s^{\prime}-\frac{d+1}{2}}} . \tag{48}
\end{equation*}
$$

Proof. The term $L_{\Omega_{\sigma}}^{\Omega_{\alpha}}=\varepsilon T_{\Omega_{\alpha}}^{\Omega_{\sigma}}$ satisfies estimate (48) by lemma 3.9. To estimate the term $L_{\Omega_{\sigma}}^{R}\left(L_{R}^{R}\right)^{-1} L_{R}^{\Omega_{\alpha}}$ we decompose $R=R_{1} \cup R_{2} \cup R_{3}$ where

$$
R_{1}:=\left\{k \in R \left\lvert\, \mathrm{d}\left(\Omega_{\alpha}, k\right)<\frac{\mathrm{d}\left(\Omega_{\alpha}, \Omega_{\sigma}\right)}{3}\right.\right\}, R_{3}:=\left\{k \in R \left\lvert\, \mathrm{d}\left(\Omega_{\sigma}, k\right)<\frac{\mathrm{d}\left(\Omega_{\alpha}, \Omega_{\sigma}\right)}{3}\right.\right\}, R_{2}:=R \backslash\left(R_{1} \cup R_{3}\right)
$$

Accordingly we have the decomposition $H_{R}=H_{R_{1}} \oplus H_{R_{2}} \oplus H_{R_{3}}$ and so

$$
L_{\Omega_{\sigma}}^{R}\left(L_{R}^{R}\right)^{-1} L_{R}^{\Omega_{\alpha}}=\sum_{i=1}^{3} \sum_{j=1}^{3} L_{\Omega_{\sigma}}^{R_{i}}\left[\left(L_{R}^{R}\right)^{-1}\right]_{R_{i}}^{R_{j}} L_{R_{j}}^{\Omega_{\alpha}} .
$$

To bound each term in the sum above, we distinguish three cases.
First case: $j \geq 2$. In this case $\mathrm{d}\left(R_{j}, \Omega_{\alpha}\right) \geq \mathrm{d}\left(\Omega_{\alpha}, \Omega_{\sigma}\right) / 3$. Then, since $L_{R_{j}}^{\Omega_{\alpha}}=\varepsilon T_{R_{j}}^{\Omega_{\alpha}}$, by lemma 3.9,

$$
\left\|L_{R_{j}}^{\Omega_{\alpha}}\right\|_{0} \leq \frac{\varepsilon C\left(s^{\prime}\right)\|a\|_{s^{\prime}}}{\mathrm{d}\left(\Omega_{\alpha}, \Omega_{\sigma}\right)^{s^{\prime}-\frac{d+1}{2}}} .
$$

Using also $\left\|L_{\Omega_{\sigma}}^{R_{i}}\right\|_{0}=\varepsilon\left\|T_{\Omega_{\sigma}}^{R_{i}}\right\|_{0} \leq \varepsilon\|a\|_{s}$ and lemma 3.7, we get

$$
\left\|L_{\Omega_{\sigma}}^{R_{i}}\left[\left(L_{R}^{R}\right)^{-1}\right]_{R_{i}}^{R_{j}} L_{R_{j}}^{\Omega_{\alpha}}\right\|_{0} \leq \frac{2 \varepsilon\|a\|_{s}}{\rho} \frac{\varepsilon C\left(s^{\prime}\right)\|a\|_{s^{\prime}}}{\mathrm{d}\left(\Omega_{\alpha}, \Omega_{\sigma}\right)^{s^{\prime}-\frac{d+1}{2}}} \leq \frac{\varepsilon C\left(s^{\prime}\right)\|a\|_{s^{\prime}}}{\mathrm{d}\left(\Omega_{\alpha}, \Omega_{\sigma}\right)^{s^{\prime}-\frac{d+1}{2}}}
$$

Second case: $j=1, i=1,2$. Now $\mathrm{d}\left(R_{i}, \Omega_{\sigma}\right) \geq \mathrm{d}\left(\Omega_{\alpha}, \Omega_{\sigma}\right) / 3$, and we proceed as in the previous case (the small factor is $\left\|L_{\Omega_{\sigma}}^{R_{i}}\right\|_{0}$ ).
Third case: $j=1, i=3$. Now $\mathrm{d}\left(R_{1}, R_{3}\right) \geq \mathrm{d}\left(\Omega_{\alpha}, \Omega_{\sigma}\right) / 3$, and, by lemma 3.11,

$$
\left\|\left[\left(L_{R}^{R}\right)^{-1}\right]_{R_{3}}^{R_{1}}\right\|_{0} \leq \frac{\varepsilon C\left(s^{\prime}\right)\|a\|_{s^{\prime}}}{\rho^{2} \mathrm{~d}\left(\Omega_{\alpha}, \Omega_{\sigma}\right)^{s^{\prime}-\frac{d+1}{2}}} .
$$

Using also $\left\|L_{\Omega_{\sigma}}^{R_{i}}\right\|_{0},\left\|L_{\Omega_{\sigma}}^{R_{i}}\right\|_{0} \leq \varepsilon\|a\|_{s}$, we get estimate (48).
In the sequel we assume

$$
\begin{equation*}
s \geq s(d, \tau):=\frac{d+1}{2}+\frac{d+1+\tau}{\delta}+1 \tag{49}
\end{equation*}
$$

where $\delta>0$ is defined in (H2) and $\tau>0$ is the parameter introduced in Proposition 3.1.

Lemma 3.13 Under the assumptions of Proposition 3.1, $\forall 0<\varepsilon \leq \varepsilon_{0}(\gamma)$ small enough, $U_{\Omega_{\alpha}}^{\Omega_{\alpha}}$ is invertible and $\left\|\left(U_{\Omega_{\alpha}}^{\Omega_{\alpha}}\right)^{-1}\right\|_{0} \leq c \gamma^{-1} M_{\alpha}^{\tau}$ with $c:=2^{\tau+2}$.
Proof It is sufficient to prove that

$$
\begin{equation*}
\left\|U_{\Omega_{\alpha}}^{\Omega_{\alpha}} w\right\|_{0} \geq \frac{\gamma}{c M_{\alpha}^{\tau}}\|w\|_{0}, \quad \forall w \in H_{\Omega_{\alpha}} \tag{50}
\end{equation*}
$$

For all $w \in H_{\Omega_{\alpha}} \subset H_{S}$, we have, recalling (47),

$$
\begin{equation*}
U_{\Omega_{\alpha}}^{\Omega_{\alpha}} w+\sum_{\sigma \neq \alpha} U_{\Omega_{\sigma}}^{\Omega_{\alpha}} w=U w=\left(L_{R}^{S}+L_{S}^{S}\right) w-\left(L_{R}^{R}+L_{S}^{R}\right)\left(L_{R}^{R}\right)^{-1} L_{R}^{S} w=\mathcal{L}_{N} h \tag{51}
\end{equation*}
$$

where $h:=w-\left(L_{R}^{R}\right)^{-1} L_{R}^{S} w$.
Step 1: $\sum_{\sigma \neq \alpha}\left\|U_{\Omega_{\sigma}}^{\Omega_{\alpha}} w\right\|_{0} \leq \varepsilon\|a\|_{s} C(s) M_{\alpha}^{d+1-\nu}\|w\|_{0}$ where $\nu:=\delta(s-(d+1) / 2)>0$.
Indeed, by lemma 3.12 (with $s^{\prime}=s$ )

$$
\begin{equation*}
\sum_{\sigma \neq \alpha}\left\|U_{\Omega_{\sigma}}^{\Omega_{\alpha}} w\right\|_{0} \leq \sum_{\sigma \neq \alpha} \frac{C \varepsilon\|a\|_{s}\|w\|_{0}}{\mathrm{~d}\left(\Omega_{\alpha}, \Omega_{\sigma}\right)^{s-\frac{d+1}{2}}} \stackrel{(H 2)}{\leq} C \varepsilon\|a\|_{s}\|w\|_{0} \sum_{\sigma \neq \alpha} \frac{1}{\left(M_{\alpha}+M_{\sigma}\right)^{\nu}} \tag{52}
\end{equation*}
$$

Next, for each $x \in \Omega_{\sigma} \subset S$, we define $M(x):=M_{\sigma}$ and $N(x):=\left|\Omega_{\sigma}\right|$. Then

$$
\sum_{\sigma \neq \alpha} \frac{1}{\left(M_{\alpha}+M_{\sigma}\right)^{\nu}}=\sum_{x \in S \backslash \Omega_{\alpha}} \frac{1}{N(x)\left(M_{\alpha}+M(x)\right)^{\nu}} \leq \sum_{x \in \mathbf{Z}^{d+1},|x| \leq N} \frac{1}{\left(M_{\alpha}+|x|\right)^{\nu}}
$$

because $N(x) \geq 1$ and $M(x) \geq|x|$. Hence, for $\nu>d+1$,

$$
\sum_{\sigma \neq \alpha} \frac{1}{\left(M_{\alpha}+M_{\sigma}\right)^{\nu}} \leq C \int_{1}^{+\infty} \frac{r^{d} d r}{\left(M_{\alpha}+r\right)^{\nu}} \leq \frac{C(\nu)}{\left(1+M_{\alpha}\right)^{\nu-d-1}}
$$

proving, by (52), Step 1 .
Step 2: $\left\|\mathcal{L}_{N} h\right\|_{0} \geq 2^{-\tau-1} \gamma M_{\alpha}^{-\tau}\|w\|_{0}$.
Decompose $h=h^{\prime}+h^{\prime \prime}$ with $h^{\prime}:=P_{\Omega_{K}} h, h^{\prime \prime}:=P_{\Omega_{K}^{c}} h$ and $K:=2 M_{\alpha}$ (recall the notation in (36)). We have

$$
\begin{equation*}
\left\|\mathcal{L}_{N} h\right\|_{0} \geq\left\|P_{\Omega_{K}} \mathcal{L}_{N} h\right\|_{0} \geq\left\|L_{K} h^{\prime}\right\|_{0}-\varepsilon\left\|T_{\Omega_{K}}^{\Omega_{K}^{c}} h^{\prime \prime}\right\|_{0} \geq \frac{\gamma}{\left(2 M_{\alpha}\right)^{\tau}}\left\|h^{\prime}\right\|_{0}-\varepsilon C\|a\|_{s}\left\|h^{\prime \prime}\right\|_{0} \tag{53}
\end{equation*}
$$

by the assumptions on the eigenvalues of $L_{K}$ in Proposition 3.1. Moreover since $h=w-\left(L_{R}^{R}\right)^{-1} L_{R}^{S} w$ and $w \in H_{\Omega_{\alpha}} \subset H_{\Omega_{K}}$,

$$
h^{\prime \prime}=-P_{\Omega_{K}^{c}}\left(L_{R}^{R}\right)^{-1} L_{R}^{S} w=-\left[\left(L_{R}^{R}\right)^{-1}\right]_{R \cap \Omega_{K}^{c}}^{R} L_{R}^{\Omega_{\alpha}} w
$$

Now $\mathrm{d}\left(\Omega_{\alpha}, R \cap \Omega_{K}^{c}\right) \geq M_{\alpha}$, from which we derive, arguing as in the proof of lemma 3.12, that

$$
\begin{equation*}
\left\|h^{\prime \prime}\right\|_{0} \leq \frac{C \varepsilon\|a\|_{s}}{\rho M_{\alpha}^{s-(d+1) / 2}}\|w\|_{0} \tag{54}
\end{equation*}
$$

Furthermore, since $w \in H_{\Omega_{\alpha}}$ and $h^{\prime}-w=-P_{\Omega_{2 M_{\alpha}}}\left(L_{R}^{R}\right)^{-1} L_{R}^{S} w \in H_{R}$, we have $\left\|h^{\prime}\right\|_{0}=\left(\left\|h^{\prime}-w\right\|_{0}^{2}+\right.$ $\left.\|w\|_{0}^{2}\right)^{1 / 2} \geq\|w\|_{0}$ and, (53), (54), imply

$$
\left\|\mathcal{L}_{N} h\right\|_{0} \geq \frac{\gamma\|w\|_{0}}{\left(2 M_{\alpha}\right)^{\tau}}-\frac{C \varepsilon^{2}\|a\|_{s}^{2}}{\rho M_{\alpha}^{s-(d+1) / 2}}\|w\|_{0} \geq \frac{\gamma\|w\|_{0}}{\left(2 M_{\alpha}\right)^{\tau}}\left(1-\frac{C^{\prime} \varepsilon^{2}\|a\|_{s}^{2}}{\gamma \rho M_{\alpha}^{s-\tau-(d+1) / 2}}\right) \geq 2^{-\tau-1} \frac{\gamma}{M_{\alpha}^{\tau}}\|w\|_{0}
$$

because $s>\tau+(d+1) / 2$, and provided that $\varepsilon \gamma^{-1}$ is small enough.
Noting that $\nu>d+1-\tau$ (because $s \geq s(d, \tau)$ ), we deduce (50) from (51), Step 1 and Step 2, for $\varepsilon_{0}(\gamma)$ small enough.

### 3.6 Inversion of $U$ and proof of Proposition 3.1

By lemma 3.13 the operator $\mathcal{D}:=\operatorname{diag}_{\alpha \in \mathcal{A}}\left(U_{\Omega_{\alpha}}^{\Omega_{\alpha}}\right)$ is invertible, and we write

$$
U=\mathcal{D}+\mathcal{R}=\mathcal{D}\left(I+\mathcal{D}^{-1} \mathcal{R}\right)
$$

For the presence of the "small divisors", $\mathcal{D}^{-1}$ acts as an unbounded operator.
Lemma 3.14 $\forall s^{\prime} \geq 0,\left\|\mathcal{D}^{-1} h\right\|_{s^{\prime}} \leq K\left(s^{\prime}\right) \gamma^{-1} N^{\tau}\|h\|_{s^{\prime}}, \forall h \in H_{\Omega_{N}}$.
Proof. For $h=\sum_{\alpha \in \mathcal{A}} h_{\alpha}, h_{\alpha} \in H_{\Omega_{\alpha}}$,

$$
\left\|\mathcal{D}^{-1} h\right\|_{s^{\prime}}^{2}=\sum_{\alpha \in \mathcal{A}}\left\|\left(U_{\Omega_{\alpha}}^{\Omega_{\alpha}}\right)^{-1} h_{\alpha}\right\|_{s^{\prime}}^{2} \stackrel{(P 4)}{\leq} \sum_{\alpha} M_{\alpha}^{2 s^{\prime}}\left\|\left(U_{\Omega_{\alpha}}^{\Omega_{\alpha}}\right)^{-1} h_{\alpha}\right\|_{0}^{2} \leq \sum_{\alpha} M_{\alpha}^{2 s^{\prime}} \frac{c M_{\alpha}^{2 \tau}}{\gamma^{2}}\left\|h_{\alpha}\right\|_{0}^{2}
$$

by lemma 3.13 . Then

$$
\left\|\mathcal{D}^{-1} h\right\|_{s^{\prime}}^{2} \stackrel{(P 4)}{\leq} \sum_{\alpha \in \mathcal{A}} \frac{c M_{\alpha}^{2 \tau}}{\gamma^{2}} M_{\alpha}^{2 s^{\prime}} \frac{\left\|h_{\alpha}\right\|_{s^{\prime}}^{2}}{m_{\alpha}^{2 s^{\prime}}} \stackrel{(H 1)}{\leq} \frac{c 2^{2 s^{\prime}}}{\gamma^{2}} \sum_{\alpha \in \mathcal{A}} M_{\alpha}^{2 \tau}\left\|h_{\alpha}\right\|_{s^{\prime}}^{2} \leq \frac{C\left(s^{\prime}\right)}{\gamma^{2}} N^{2 \tau}\|h\|_{s^{\prime}}^{2}
$$

because $M_{\alpha} \leq N$.
By lemma 3.12 the operator $\mathcal{R}$ acts somehow as a multiplication operator. Nevertheless, using the separation property (H2) we prove that $\mathcal{D}^{-1} \mathcal{R}$ is $L^{2}$-bounded and, using also the dyadic property (H1), we prove an interpolation type estimate for $\mathcal{D}^{-1} \mathcal{R}$ in high Sobolev norm.

Lemma 3.15 Assume (49). Then

$$
\begin{equation*}
\left\|\mathcal{D}^{-1} \mathcal{R} h\right\|_{0} \leq K \frac{\varepsilon}{\gamma}\|a\|_{s}\|h\|_{0} \tag{55}
\end{equation*}
$$

and, $\forall s^{\prime} \geq s(d, \tau)$, setting $\mu_{0}:=\tau+3(d+1) / 2>0$,

$$
\begin{equation*}
\left\|\mathcal{D}^{-1} \mathcal{R} h\right\|_{s^{\prime}} \leq K\left(s^{\prime}\right) \frac{\varepsilon}{\gamma}\left(\|h\|_{s^{\prime}}\|a\|_{s}+\|a\|_{s^{\prime}}\|h\|_{0} N^{\mu_{0}}\right) . \tag{56}
\end{equation*}
$$

Proof. We shall use several times the following bound, proved as in the first step of the proof of Lemma 3.13: if $\lambda>0$ and $\eta>d+1$, then

$$
\begin{equation*}
\sum_{\sigma \in \mathcal{A}} \frac{1}{\left(\lambda+M_{\sigma}\right)^{\eta}} \leq \frac{C(\eta)}{(1+\lambda)^{\eta-(d+1)}} . \tag{57}
\end{equation*}
$$

Given $h=\sum_{\alpha \in \mathcal{A}} h_{\alpha}$, we have

$$
\mathcal{D}^{-1} \mathcal{R} h=\sum_{\sigma \in \mathcal{A}} w_{\sigma} \quad \text { with } \quad w_{\sigma}:=\left(U_{\Omega_{\sigma}}^{\Omega_{\sigma}}\right)^{-1}\left(\sum_{\alpha \neq \sigma} U_{\Omega_{\sigma}}^{\Omega_{\alpha}} h_{\alpha}\right) \in H_{\Omega_{\sigma}}
$$

From lemmas 3.13, 3.12 and (H2), we get, for $\nu:=\delta(s-(d+1) / 2)$,

$$
\begin{align*}
\left\|w_{\sigma}\right\|_{0} & \leq c \frac{M_{\sigma}^{\tau}}{\gamma}\left(\sum_{\alpha \neq \sigma}\left\|U_{\Omega_{\sigma}}^{\Omega_{\alpha}}\right\|_{0}\left\|h_{\alpha}\right\|_{0}\right)  \tag{58}\\
& \leq c \frac{M_{\sigma}^{\tau}}{\gamma}\left(\sum_{\alpha \neq \sigma} \frac{C \varepsilon\|a\|_{s}}{\left(M_{\alpha}+M_{\sigma}\right)^{\nu}}\left\|h_{\alpha}\right\|_{0}\right) \\
& \leq C \frac{\varepsilon}{\gamma}\|a\|_{s}\|h\|_{0} M_{\sigma}^{\tau}\left(\sum_{\alpha \neq \sigma} \frac{1}{\left(M_{\sigma}+M_{\alpha}\right)^{2 \nu}}\right)^{1 / 2} \leq \frac{\varepsilon}{\gamma} \frac{C\|a\|_{s}}{\left(1+M_{\sigma}\right)^{\nu-\tau-(d+1) / 2}}\|h\|_{0},
\end{align*}
$$

by the Cauchy-Schwartz inequality and (57). Hence

$$
\left\|\mathcal{D}^{-1} \mathcal{R} h\right\|_{0}=\left(\sum_{\sigma \in \mathcal{A}}\left\|w_{\sigma}\right\|_{0}^{2}\right)^{1 / 2} \leq C \frac{\varepsilon}{\gamma}\left(\sum_{\sigma} \frac{1}{\left(1+M_{\sigma}\right)^{2 \nu-2 \tau-(d+1)}}\right)^{1 / 2}\|a\|_{s}\|h\|_{0} \leq K \frac{\varepsilon}{\gamma}\|a\|_{s}\|h\|_{0}
$$

still by (57) (note that $\nu>\tau+d+1$ by (49)). This proves (55).
In order to prove (56) we observe that:
i) if $M_{\alpha} \geq M_{\sigma} / 4$ then

$$
M_{\sigma}^{s^{\prime}}\left\|U_{\Omega_{\sigma}}^{\Omega_{\alpha}}\right\|_{0}\left\|h_{\alpha}\right\|_{0} \stackrel{(P 4)}{\leq}\left\|U_{\Omega_{\sigma}}^{\Omega_{\alpha}}\right\|_{0} \frac{M_{\sigma}^{s^{\prime}}}{m_{\alpha}^{s^{\prime}}}\left\|h_{\alpha}\right\|_{s^{\prime}} \stackrel{(H 1)}{\leq}\left\|U_{\Omega_{\sigma}}^{\Omega_{\alpha}}\right\|_{0} 2^{s^{\prime}} \frac{M_{\sigma}^{s^{\prime}}}{M_{\alpha}^{s^{\prime}}}\left\|h_{\alpha}\right\|_{s^{\prime}} \leq \frac{C \varepsilon 8^{s^{\prime}}\|a\|_{s}}{\left(M_{\alpha}+M_{\sigma}\right)^{)^{\prime}}}\left\|h_{\alpha}\right\|_{s^{\prime}},
$$

using also lemma 3.12 and (H2).
ii) if $M_{\alpha}<M_{\sigma} / 4$ then, by (H1), $\operatorname{dist}\left(\Omega_{\alpha}, \Omega_{\sigma}\right) \geq M_{\sigma} / 4$, and, by Lemma 3.12,

$$
M_{\sigma}^{s^{\prime}}\left\|U_{\Omega_{\sigma}}^{\Omega_{\alpha}}\right\|_{0}\left\|h_{\alpha}\right\|_{0} \leq M_{\sigma}^{s^{\prime}} \frac{\varepsilon C\left(s^{\prime}\right)\|a\|_{s^{\prime}}}{\mathrm{d}\left(\Omega_{\alpha}, \Omega_{\sigma}\right)^{s^{\prime}-\frac{d+1}{2}}}\left\|h_{\alpha}\right\|_{0} \leq K\left(s^{\prime}\right)\|a\|_{s^{\prime}} \varepsilon M_{\sigma}^{(d+1) / 2}\left\|h_{\alpha}\right\|_{0}
$$

Therefore, from (58),

$$
\begin{align*}
\left\|w_{\sigma}\right\|_{s^{\prime}} & \stackrel{(P 4)}{\leq} M_{\sigma}^{s^{\prime}}\left\|w_{\sigma}\right\|_{0} \leq C\left(s^{\prime}\right) \frac{\varepsilon}{\gamma} M_{\sigma}^{\tau}\left(\sum_{\alpha, M_{\alpha} \geq M_{\sigma} / 4} \frac{\|a\|_{s}\left\|h_{\alpha}\right\|_{s^{\prime}}}{\left(M_{\alpha}+M_{\sigma}\right)^{\nu}}+\|a\|_{s^{\prime}} M_{\sigma}^{(d+1) / 2} \sum_{\alpha, M_{\alpha}<M_{\sigma} / 4}\left\|h_{\alpha}\right\|_{0}\right) \\
& \leq C\left(s^{\prime}\right) \frac{\varepsilon}{\gamma} M_{\sigma}^{\tau}\left(\frac{\|a\|_{s}\|h\|_{s^{\prime}}}{M_{\sigma}^{\nu-(d+1) / 2}}+\|a\|_{s^{\prime}} M_{\sigma}^{d+1}\|h\|_{0}\right) \tag{59}
\end{align*}
$$

by the Cauchy-Schwartz inequality and (57). Finally

$$
\begin{aligned}
\left\|\mathcal{D}^{-1} \mathcal{R} h\right\|_{s^{\prime}} & =\left(\sum_{\sigma \in \mathcal{A}}\left\|w_{\sigma}\right\|_{s^{\prime}}^{2}\right)^{1 / 2} \\
& \stackrel{(59)}{\leq} C\left(s^{\prime}\right) \frac{\varepsilon}{\gamma}\left[\left(\sum_{\sigma} \frac{1}{M_{\sigma}^{2 \nu-2 \tau-(d+1)}}\right)^{1 / 2}\|a\|_{s}\|h\|_{s^{\prime}}+\left(\sum_{\sigma} M_{\sigma}^{2 \tau+2(d+1)}\right)^{1 / 2}\|a\|_{s^{\prime}}\|h\|_{0}\right] \\
& \stackrel{(57)}{\leq} K\left(s^{\prime}\right) \frac{\varepsilon}{\gamma}\left(\|h\|_{s^{\prime}}\|a\|_{s}+\|a\|_{s^{\prime}} N^{\tau+3(d+1) / 2}\|h\|_{0}\right)
\end{aligned}
$$

proving (56).
By (55), for $\varepsilon \gamma^{-1}\|a\|_{s}$ small, the operator $U$ is invertible, and, using (56), we prove the following interpolation type inequality for $U^{-1}$.

Lemma 3.16 $\forall s^{\prime} \geq s(d, \tau)$, there is $c\left(s^{\prime}\right)>0$ such that, if $\varepsilon \gamma^{-1}\left(\|a\|_{s}+1\right) \leq c\left(s^{\prime}\right)$, then

$$
\left\|U^{-1} h\right\|_{s^{\prime}} \leq K\left(s^{\prime}\right) \gamma^{-1} N^{\mu}\left(\|h\|_{s^{\prime}}+\|a\|_{s^{\prime}}\|h\|_{0}\right), \quad \forall h \in H_{S},
$$

with $\mu:=2 \tau+3(d+1) / 2$.
Proof. If $\varepsilon\|a\|_{s} \gamma^{-1}$ is small enough then, by (55), $\left\|D^{-1} \mathcal{R}\right\|_{0} \leq 1 / 2$ and $I+D^{-1} \mathcal{R}$ is invertible, $\left\|\left(I+D^{-1} \mathcal{R}\right)^{-1}\right\|_{0} \leq 2$. Then, arguing like in the proof of Lemma 3.8, we derive from (56) that for $\varepsilon\left(\|a\|_{s}+1\right) \gamma^{-1} \leq c\left(s^{\prime}\right)$ small enough,

$$
\left\|\left(I+\mathcal{D}^{-1} \mathcal{R}\right)^{-1} h\right\|_{s^{\prime}} \leq 2\|h\|_{s^{\prime}}+\|a\|_{s^{\prime}}\|h\|_{0} N^{\mu_{0}}
$$

Using lemma 3.14 the thesis follows.
Proof of Proposition 3.1 concluded. Lemma 3.16, (46), (45) and lemma 3.8 yield (37). We use also that $\|a\|_{s^{\prime}}=\left\|\left(\partial_{u} F\right)(t, x, u)\right\|_{s^{\prime}} \leq C\left(s^{\prime}\right)\|u\|_{s^{\prime}}$, by property (P2) applied to $\partial_{u} F$.

## 4 Separation properties of the singular sites

In this section we verify, assuming only the standard Diophantine condition (7) and the nonresonance condition (2), that, for $\rho<\rho(\gamma)$, there is a partition of the singular sites like in (35) verifying (H1)-(H2) with $\delta$ depending only on $d$. The proof follows essentially the scheme of [6] except in lemmas 4.1 and 4.4.

Consider the bilinear symmetric form $\varphi_{\omega}: \mathbf{R}^{d+1} \times \mathbf{R}^{d+1} \rightarrow \mathbf{R}$ defined by

$$
\varphi_{\omega}\left(x, x^{\prime}\right):=j \cdot j^{\prime}-\omega^{2} l l^{\prime}, \quad \forall x=(l, j), x^{\prime}=\left(l^{\prime}, j^{\prime}\right) \in \mathbf{R} \times \mathbf{R}^{d}
$$

and the corresponding quadratic form

$$
Q_{\omega}(x)=\varphi_{\omega}(x, x):=|j|^{2}-\omega^{2} l^{2} .
$$

A vector $x=(l, j) \in \mathbf{Z} \times \mathbf{Z}^{d}$ is said "weakly singular" if $\left|Q_{\omega}(x)\right| \leq 1+|m|$.
Definition 4.1 $A$ sequence $x_{0}, \ldots, x_{K} \in \mathbf{Z}^{d+1}$ of distinct, weakly singular, integer vectors satisfying, for some $B \geq 2,\left|x_{k+1}-x_{k}\right| \leq B, \forall k=0, \ldots, K-1$, is called a $B$-chain of length $K$.

Theorem 4.1 If $\omega^{2}$ satisfies (7), then any $B$-chain has length $K \leq B^{C} / \gamma^{r}$ for some $C:=C(d)>0$ and $r:=r(d)>0$.

Remark 4.1 Theorem 4.1 has been proved in [6], see lemma 2.10, assuming the stronger Diophantine condition $\left|\sum_{j=0}^{10 d} a_{j} \omega^{2 j}\right| \geq\left(\sum_{j}\left|a_{j}\right|\right)^{-C}, \forall\left(a_{j}\right) \in \mathbf{Z}^{10 d+1} \backslash\{0\}$ and in [10], lemma 8.7, assuming $\left|\sum_{j=0}^{d+1} a_{j} \omega^{2 j}\right| \geq\left(\sum_{j}\left|a_{j}\right|\right)^{-C}, \forall\left(a_{j}\right) \in \mathbf{Z}^{d+1} \backslash\{0\}$.

The proof of Theorem 4.1 is split in several lemmas.
Given integer vectors $f_{i} \in \mathbf{Z}^{d+1}, i=1, \ldots, n, 1 \leq n \leq d+1$, linearly independent on $\mathbf{R}$, we consider the subspace $F:=\operatorname{Span}_{\mathbf{R}}\left\{f_{1}, \ldots, f_{n}\right\}$ of $\mathbf{R}^{d+1}$ and the restriction $\varphi_{\omega \mid F}$ of the bilinear form $\varphi_{\omega}$ to $F$, which is represented by the symmetric matrix

$$
A_{\omega}:=\left\{\varphi_{\omega}\left(f_{i}, f_{i^{\prime}}\right)\right\}_{i, i^{\prime}=1}^{n} \in \operatorname{Mat}_{n}(\mathbf{R}) .
$$

Introducing the symmetric bilinear forms

$$
\mathcal{R}\left(x, x^{\prime}\right):=j \cdot j^{\prime} \quad \text { and } \quad \mathcal{S}\left(x, x^{\prime}\right):=l l^{\prime}
$$

we write

$$
\varphi_{\omega}=\mathcal{R}-\omega^{2} \mathcal{S} \text { and } A_{\omega}=R-\omega^{2} S
$$

where $R:=\left\{\mathcal{R}\left(f_{i}, f_{i^{\prime}}\right)\right\}_{i, i^{\prime}=1}^{n}=\left(R_{1}, \ldots, R_{n}\right), S:=\left\{\mathcal{S}\left(f_{i}, f_{i^{\prime}}\right)\right\}_{i, i^{\prime}=1}^{n}=\left(S_{1}, \ldots, S_{n}\right)$ are the matrices that represent respectively $\mathcal{R}_{\mid F}$ and $\mathcal{S}_{\mid F}$ in the basis $\left\{f_{1}, \ldots, f_{n}\right\}$. Note that the matrices $R, S$ have integer coefficients. Here $R_{i}, S_{i} \in \mathbf{Z}^{n}, i=1, \ldots, n$ denote the column vectors respectively of $R$ and $S$. The following lemma is the main difference with respect to [6].

Lemma 4.1 Assume that $\omega^{2}$ satisfies (7). Then $A_{\omega}$ is invertible and

$$
\begin{equation*}
\left\|A_{\omega}^{-1}\right\| \leq \frac{c(n)}{\gamma}\left(\max _{i=1, \ldots, n}\left|f_{i}\right|\right)^{5 n-2} \tag{60}
\end{equation*}
$$

Proof. The matrix $S$ has rank at most 1 because it represents the restriction to $F$ of a bilinear form of rank 1. Since any two columns of $S$ are colinear, the development in $\omega^{2}$ of $\operatorname{det} A_{\omega}$ reduces to

$$
\begin{equation*}
\operatorname{det} A_{\omega}=\operatorname{det}\left(R_{1}-\omega^{2} S_{1}, \ldots, R_{n}-\omega^{2} S_{n}\right)=\operatorname{det}\left(R_{1}, \ldots, R_{n}\right)-\omega^{2} \sum_{i=1}^{n} \operatorname{det}\left(R_{1}, \ldots, S_{i}, \ldots, R_{n}\right) \tag{61}
\end{equation*}
$$

Therefore $\operatorname{det} A_{\omega}=P\left(\omega^{2}\right)$ is a polynomial in $\omega^{2}$ of degree at most 1 , with integer coefficients since det $R$, $\operatorname{det}\left(R_{1}, \ldots, S_{i}, \ldots, R_{n}\right) \in \mathbf{Z}$. Furthermore $P(\cdot)$ is not identically zero because $P(-1)=\operatorname{det}(R+S)$ is positive (the matrix $R+S$ being positive definite).

By (61), if $\operatorname{det} R=0$ then $\left|\operatorname{det} A_{\omega}\right| \geq \omega^{2}$, and, if $\operatorname{det} R \neq 0$, since $\omega^{2}$ satisfies (7), $\left|\operatorname{det} A_{\omega}\right| \geq$ $\gamma|\operatorname{det} R|^{-3 / 2}$. Since $R+S={ }^{t} \mathcal{F} \mathcal{F}$ with $\mathcal{F}=\left(f_{1}, \ldots, f_{n}\right)$, we have $0 \leq \operatorname{det} R \leq \operatorname{det}(R+S)=(\operatorname{det} \mathcal{F})^{2} \leq$ $\left|f_{1}\right|^{2} \ldots\left|f_{n}\right|^{2} \leq M^{2 n}$, where $M:=\max _{i=1, \ldots, n}\left|f_{i}\right|$. Hence

$$
\begin{equation*}
\left|\operatorname{det} A_{\omega}\right| \geq \frac{\gamma}{M^{3 n}} \tag{62}
\end{equation*}
$$

By (7), $\omega^{2} \geq \gamma$, and (62) holds in both cases. Applying the Cramer rule

$$
\left|\left(A_{\omega}^{-1}\right)_{i, i^{\prime}}\right| \leq \frac{c(n)}{\left|\operatorname{det} A_{\omega}\right|}\left(M^{2}\right)^{(n-1)} \stackrel{(62)}{\leq} \frac{c^{\prime}(n)}{\gamma} M^{3 n+2(n-1)}
$$

whence (60) follows.
The remainder of the proof follows [6] (although some details are different). We report it for completeness.

By lemma 4.1 the bilinear form $\varphi_{\omega \mid F}$ is non-degenerate and we decompose

$$
\mathbf{R}^{d+1}=F \oplus F^{\perp \varphi_{\omega}} \quad \text { where } \quad F^{\perp \varphi_{\omega}}:=\left\{x \in \mathbf{R}^{d+1} \mid \varphi_{\omega}(x, f)=0, \quad \forall f \in F\right\}
$$

We denote by $P_{F}: \mathbf{R}^{d+1} \rightarrow F$ the corresponding projector onto $F$.
Lemma 4.2 Let $x_{0}, \ldots, x_{L}$ be a $B$-chain. Define $G:=\operatorname{Span}_{\mathbf{R}}\left\{x_{i}-x_{k} \mid i, k=0, \ldots, L\right\} \subset \mathbf{R}^{d+1}$. Then, calling $n:=\operatorname{dim} G$,

$$
\begin{equation*}
\left|P_{G}\left(x_{j}\right)\right| \leq \frac{c(n)}{\gamma}(L B)^{5 n+1}, \quad \forall j=0, \ldots, L \tag{63}
\end{equation*}
$$

Proof. For all $i \neq j$, by bilinearity,

$$
\begin{equation*}
Q_{\omega}\left(x_{i}\right)=Q_{\omega}\left(x_{j}+\left(x_{i}-x_{j}\right)\right)=Q_{\omega}\left(x_{j}\right)+2 \varphi_{\omega}\left(x_{j}, x_{i}-x_{j}\right)+Q_{\omega}\left(x_{i}-x_{j}\right) . \tag{64}
\end{equation*}
$$

Since $\left|x_{i}-x_{j}\right| \leq|i-j| B$ and $\left|Q_{\omega}\left(x_{i}\right)\right|,\left|Q_{\omega}\left(x_{j}\right)\right| \leq 1+|m|$, we get

$$
\begin{equation*}
\left|2 \varphi_{\omega}\left(x_{j}, x_{i}-x_{j}\right)\right| \stackrel{(64)}{\leq} 2+2|m|+C\left|x_{i}-x_{j}\right|^{2} \leq 2+2|m|+C|i-j|^{2} B^{2} \leq C^{\prime} L^{2} B^{2} \tag{65}
\end{equation*}
$$

According to the decomposition $\mathbf{R}^{d+1}=G \oplus G^{\perp \varphi_{\omega}}$ we write $x_{j}=y_{j}+z_{j}$ with $y_{j}:=P_{G}\left(x_{j}\right) \in G$ and $z_{j} \in G^{\perp \varphi_{\omega}}$.

Fixed any $j \in[0, L]$ we have $G=\operatorname{Span}_{\mathbf{R}}\left\{x_{i}-x_{j} ; i=0, \ldots, L\right\}$. Let us extract from $\left\{x_{i}-x_{j}\right.$, $i=0, \ldots, L\}$, a basis $\left\{f_{1}, \ldots, f_{n}\right\}$ of $G$. We develop $y_{j}=\sum_{l=1}^{n} a_{l, j} f_{l}$ whence

$$
A_{\omega} a=b \quad \text { where } \quad a:=\left(\begin{array}{c}
a_{1, j}  \tag{66}\\
\ldots \\
a_{n, j}
\end{array}\right), \quad b:=\left(\begin{array}{c}
\varphi_{\omega}\left(f_{1}, x_{j}\right) \\
\ldots \\
\varphi_{\omega}\left(f_{n}, x_{j}\right)
\end{array}\right)
$$

using that $\varphi_{\omega}\left(f_{m}, y_{j}\right)=\varphi_{\omega}\left(f_{m}, x_{j}\right)$ because $z_{j} \in G^{\perp \varphi_{\omega}}$. By (65), $|b| \leq C L^{2} B^{2}$ and, by lemma 4.1,

$$
|a| \leq\left\|A_{\omega}^{-1}\right\||b| \leq \frac{c(n)}{\gamma}\left(\max _{i=1, \ldots, n}\left|f_{i}\right|\right)^{5 n-2} L^{2} B^{2} \leq \frac{c(n)}{\gamma}(L B)^{5 n}
$$

since $\left|f_{i}\right| \leq L B$. Finally, $\left|y_{j}\right| \leq \sum_{l=1}^{n}\left|a_{l, j}\right|\left|f_{l}\right| \leq(c(n) / \gamma)(L B)^{5 n} L B$, which gives (63).
Lemma 4.3 Let $x_{0}, \ldots, x_{K}$ be a $B$-chain. Let $1 \leq L \leq K$ be such that

$$
\begin{equation*}
\forall j=0, \ldots, K, \quad \operatorname{Span}_{\mathbf{R}}\left\{x_{i}-x_{j}| | i-j \mid \leq L\right\}=\operatorname{Span}_{\mathbf{R}}\left\{x_{i}-x_{l} \mid i, l=0, \ldots, K\right\}=: F \tag{67}
\end{equation*}
$$

Then, for $\rho_{n}:=5 n+1, n:=\operatorname{dim} F$,

- a) $\left|P_{F}\left(x_{j}\right)\right| \leq \frac{c(n)}{\gamma}(L B)^{\rho_{n}}, \forall j=0, \ldots, K$;
- b) $\left|x_{j_{2}}-x_{j_{1}}\right| \leq \frac{c(n)}{\gamma}(L B)^{\rho_{n}}, \forall j_{1}, j_{2} \in\{0, \ldots, K\}$;
- c) $K \leq \frac{c(n)}{\gamma^{n}}(L B)^{\rho_{n} n}$.

Proof. For every fixed $j \in\{0, \ldots, K\}$ apply lemma 4.2 to the $B$-chain $\left\{x_{i}\right\}$ with $|i-j| \leq L, i \in$ $\{0, \ldots, K\}$. By (67) we have $G:=\operatorname{Span}_{\mathbf{R}}\left\{x_{i}-x_{k}| | i-j|,|k-j| \leq L\}=F\right.$, and therefore, by (63) we deduce $a$ ).

Now, by the definition of $F, \forall j_{1}, j_{2} \in\{0, \ldots, K\}$, we have $x_{j_{2}}-x_{j_{1}} \in F$ and therefore by $\left.a\right)$,

$$
\left|x_{j_{2}}-x_{j_{1}}\right|=\left|P_{F}\left(x_{j_{2}}-x_{j_{1}}\right)\right| \leq\left|P_{F}\left(x_{j_{2}}\right)\right|+\left|P_{F}\left(x_{j_{1}}\right)\right| \leq \frac{c(n)}{\gamma}(L B)^{5 n+1}
$$

which is $b$ ).
Finally, $\forall j=1, \ldots, K$, by $b),\left|x_{j}-x_{0}\right| \leq(c(n) / \gamma)(L B)^{\rho_{n}}$, whence $x_{j}$ belongs to the $n$-dimensional ball centered at $x_{0}$ with radius $(C(n) / \gamma)(L B)^{\rho_{n}}$. The number of integer vectors inside such a ball is less or equal to $\left(C^{\prime}(n) / \gamma^{n}\right)(L B)^{\rho_{n} n}$ and we deduce $c$ ).
Proof of Theorem 4.1 concluded. Define $F:=\operatorname{Span}_{\mathbf{R}}\left\{x_{i}-x_{l} \mid i, l=0, \ldots, K\right\} \subset \mathbf{R}^{d+1}$.
If $\operatorname{dim} F=1$ then the $B$-chain $x_{0}, \ldots, x_{K}$ satisfies (67) with $L=1$ because all the $x_{i}$ are distinct. By $c$ ) we get $K \leq(c / \gamma) B^{6}$. Next, suppose by induction that Theorem 4.1 holds if $\operatorname{dim} F \leq n, n \leq d$. We want to prove Theorem 4.1 when $\operatorname{dim} F=n+1$.

Fix $L=\left[K^{1 / \alpha}\right]$ with $\alpha:=2 \rho_{n+1}(n+1)$. If the $B$-chain $x_{0}, \ldots, x_{K}$ satisfies (67), then, by $\left.c\right)$,

$$
K \leq \frac{c(n+1)}{\gamma^{n+1}}(L B)^{\rho_{n+1}(n+1)} \leq \frac{c(n+1)}{\gamma^{n+1}} \sqrt{K} B^{\rho_{n+1}(n+1)} \quad \Longrightarrow \quad K \leq \frac{c^{\prime}(n+1)}{\gamma^{2(n+1)}} B^{2 \rho_{n+1}(n+1)} \leq \frac{B^{C(n)}}{\gamma^{r(n)}} .
$$

Otherwise, there exists $j \in\{0, \ldots, K\}$ such that $\operatorname{dim} \operatorname{Span}_{\mathbf{R}}\left\{x_{i}-x_{j}| | i-j \mid \leq L\right\} \leq n$. Consider the $B$-chain $\left\{x_{i}\right\}$ with $|i-j| \leq L$ whose length is at most $2 L$. By the inductive assumption, $2 L \leq$ $B^{C(n-1)} / \gamma^{r(n-1)}$, whence $K \leq(L+1)^{\alpha} \leq c B^{\alpha C(n-1)} / \gamma^{\alpha r(n-1)} \leq B^{C(n)} / \gamma^{r(n)}$. In both cases $K \leq B^{C} / \gamma^{r}$ for some $C:=C(n), r=r(n)$.

Finally we show how to get the required decomposition of the singular sites in clusters.
Lemma 4.4 Assume (2) and (7). $\forall \gamma \in(0,1)$, if $\rho \leq \rho(\gamma):=\gamma^{3 r(d)+1} / 8$ there exists a decomposition of the singular sites

$$
\begin{equation*}
S:=\left\{(l, j) \in \Omega_{N}:\left|\omega^{2} l^{2}-\lambda_{j}\right|<\rho\right\}=\bigcup_{\alpha \in \mathcal{A}} \Omega_{\alpha} \tag{68}
\end{equation*}
$$

in clusters, like in (35), satisfying properties (H1)-(H2) with $\delta:=1 / 2(C(d)+1)$ (the constants $C(d), r(d)$ are those of Theorem 4.1).

Proof. We introduce the following equivalence relation.
Definition 4.2 Given $\delta>0$, two integer vectors $x, y \in S$ are said equivalent if there exist $x_{l} \in S$, $l=0,1, \ldots n$, with $x_{0}=x, x_{n}=y$ and $\left|x_{l+1}-x_{l}\right| \leq\left(\left|x_{l}\right|+\left|x_{l+1}\right|\right)^{\delta}, \forall l$.

Chosen $\delta:=\delta(d):=1 / 2(C(d)+1)$, where $C(d)$ is the constant which appears in Theorem 4.1, we get a decomposition of $S$ in disjoint equivalence classes $\Omega_{\alpha}$ as in (68). To verify that each $\Omega_{\alpha}$ is dyadic, consider $z_{\alpha} \in \Omega_{\alpha}$ such that $\left|z_{\alpha}\right|=\max _{z \in \Omega_{\alpha}}|z|=M_{\alpha}$. Each $x \in \Omega_{\alpha}$ is connected to $z_{\alpha}$ by a $B$-chain with $B=\left(2 M_{\alpha}\right)^{\delta}$. Hence, by Theorem 4.1, there are no more than $B^{C} / \gamma^{r}$ elements in this chain. Therefore $\forall x \in \Omega_{\alpha}$,

$$
|x| \geq\left|z_{\alpha}\right|-\frac{B^{C} B}{\gamma^{r}}=M_{\alpha}-\frac{\left(2 M_{\alpha}\right)^{\delta(C+1)}}{\gamma^{r}}=M_{\alpha}-\frac{\sqrt{2 M_{\alpha}}}{\gamma^{r}} \geq \frac{M_{\alpha}}{2},
$$

provided that $M_{\alpha} \geq 8 / \gamma^{2 r}$. As a consequence, $\Omega_{\alpha}$ satisfies the dyadic property (H1) whenever $M_{\alpha} \geq$ $8 / \gamma^{2 r}$. Now, if $|x|=|(l, j)| \leq 8 / \gamma^{2 r}$ then by (2)

$$
\left|\omega^{2} l^{2}-\lambda_{j}\right| \geq \frac{3 \gamma}{\left(8 \gamma^{-2 r}\right)^{3 / 2}} \geq \frac{\gamma^{3 r+1}}{8}
$$

As a result, if $\rho \leq \rho(\gamma):=\gamma^{3 r+1} / 8$ then there is no singular site in the open ball of center 0 and radius $8 / \gamma^{2 r}$, and all the $\Omega_{\alpha}$ satisfy $M_{\alpha} \geq 8 / \gamma^{2 r}$ and are dyadic.

To prove (H2) consider $x_{\alpha} \in \Omega_{\alpha}, x_{\sigma} \in \Omega_{\sigma}$ such that $\mathrm{d}\left(\Omega_{\alpha}, \Omega_{\sigma}\right)=\left|x_{\alpha}-x_{\sigma}\right|$. Since $x_{\sigma} \notin \Omega_{\alpha}$

$$
\left|x_{\alpha}-x_{\sigma}\right|>\left(\left|x_{\alpha}\right|+\left|x_{\sigma}\right|\right)^{\delta} \geq\left(m_{\alpha}+m_{\sigma}\right)^{\delta} \geq \frac{1}{2}\left(M_{\alpha}+M_{\sigma}\right)^{\delta}
$$

by (H1), that is (H2).

## References

[1] Baldi P., Periodic solutions of forced Kirkhoff equation, preprint 2007.
[2] Baldi P., Berti M., Forced vibrations of a nonhomogeneous string, to appear in SIAM J. Math. Anal.
[3] Berti M., Bolle P., Cantor families of periodic solutions for completely resonant nonlinear wave equations, Duke Math. J. 134 (2006) 359-419.
[4] Berti M., Bolle P., Cantor families of periodic solutions of wave equations with $C^{k}$ nonlinearities to appear in NoDEA.
[5] Brezis, H.; Coron, J.-M.; Nirenberg, L., Free vibrations for a nonlinear wave equation and a theorem of P. Rabinowitz, Comm. Pure Appl. Math. 33 (1980) 667-684.
[6] Bourgain J., Construction of periodic solutions of nonlinear wave equations in higher dimension, Geom. Funct. Anal. 5 (1995) 629-639.
[7] Bourgain J., Quasi-periodic solutions of Hamiltonian perturbations of $2 D$ linear Schrödinger equations, Ann. of Math. 148 (1998) 363-439.
[8] Bourgain J., Green's function estimates for lattice Schrödinger operators and applications, Annals of Mathematics Studies 158, Princeton University Press, Princeton, 2005
[9] Chierchia L., You J., KAM tori for 1D nonlinear wave equations with periodic boundary conditions, Comm. Math. Phys. 211 (2000) 497-525.
[10] Craig W., Problèmes de petits diviseurs dans les équations aux dérivées partielles, Panoramas et Synthèses, 9, Société Mathématique de France, Paris, 2000
[11] Craig W., Wayne C. E., Newton's method and periodic solutions of nonlinear wave equation, Comm. Pure Appl. Math. 46 (1993) 1409-1498.
[12] Eliasson L. H., Kuksin S., KAM for the nonlinear Schrödinger equation, to appear in Annals of Math..
[13] Fokam J. M. , Forced vibrations via Nash-Moser iteration, PhD thesis at Austin University.
[14] Fröhlich J., Spencer T., Absence of diffusion in the Anderson tight binding model for large disorder or low energy, Comm. Math. Phys. 88 (1983) 151-184.
[15] Gentile G., Procesi M., Periodic solutions for the Schrödinger equation with nonlocal smoothing nonlinearities in higher dimension, to appear in Journal of PDEs.
[16] Kuksin S., Hamiltonian perturbations of infinite-dimensional linear systems with imaginary spectrum, Funktsional Anal. i Prilozhen. 21 (1987) 22-37, 95.
[17] Kuksin S., Analysis of Hamiltonian PDEs, Oxford Lecture series in Mathematics and its applications 19, Oxford University Press, 2000.
[18] Kuksin S., Pöschel J., Invariant Cantor manifolds of quasi-periodic oscillations for a nonlinear Schrödinger equation, Ann. of Math. (2) 143 (1996) 149-179.
[19] Lojasiewicz S., Zehnder E., An inverse function theorem in Fréchet-spaces, J. Funct. Anal. 33 (1979) 165-174.
[20] Moser J., A rapidly convergent iteration method and non-linear partial differential equations I \& II, Ann. Scuola Norm. Sup. Pisa (3) 20 (1966) 265-315 \& 499-535.
[21] Plotnikov P.I., Yungerman L.N., Periodic solutions of a weakly nonlinear wave equation with an irrational relation of period to interval length, translation in Diff. Equations 24 (1988), no. 9, 10591065 (1989).
[22] Pöschel J., A KAM-Theorem for some nonlinear partial differential equations, Ann. Scuola Norm. Sup. Pisa Cl. Sci.(4) 23 (1996) 119-148.
[23] Rabinowitz P., Free vibrations for a semi-linear wave equation, Comm. Pure Appl. Math. 31 (1978) 31-68.
[24] Taylor M., Partial Differential equations, Vol. III, Springer-Verlag, New York, 1997.
[25] Wayne E., Periodic and quasi-periodic solutions of nonlinear wave equations via KAM theory, Comm. Math. Phys. 127 (1990) 479-528.


[^0]:    *Dipartimento di Matematica e Applicazioni R. Caccioppoli Università degli Studi Napoli Federico II, Via Cintia, Monte S. Angelo I-80126 Napoli, Italy, m.berti@unina.it.
    ${ }^{\dagger}$ Université d'Avignon et des Pays de Vaucluse, Laboratoire d'Analyse non Linéaire et Géométrie (EA 2151), F-84018 Avignon, France, philippe.bolle@univ-avignon.fr.
    ${ }^{1}$ Supported by MIUR Variational Methods and Nonlinear Differential Equations.

[^1]:    ${ }^{2}$ That is $F(\cdot, x, u)$ is $2 \pi$-periodic.

[^2]:    ${ }^{3}$ The symbol $z^{*}$ denotes the complex conjugate of $z \in \mathbf{C}$.

[^3]:    ${ }^{4}$ The symbol [ $\cdot$ ] denotes the integer part.

[^4]:    ${ }^{5}$ Comparing with the notations of section $2, H_{\Omega_{N}} \equiv W^{(N)}, D \equiv L_{\omega}$ and $P_{\Omega_{N}} \equiv P_{N}$.

