# Forced vibrations of a nonhomogeneous string 

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#### Abstract

We prove existence of vibrations of a nonhomogeneous string under a nonlinear time periodic forcing term, in the case the forcing frequency avoids resonances with the vibration modes of the string (non-resonant case). The proof relies on a Nash-Moser iteration scheme. ${ }^{1}$


## 1 Introduction

In this paper we study forced vibrations of a nonhomogeneous string

$$
\left\{\begin{array}{l}
\rho(x) u_{t t}-\left(p(x) u_{x}\right)_{x}=\mu f(x, t, u)  \tag{1}\\
u(0, t)=u(\pi, t)=0
\end{array}\right.
$$

where $\rho(x)>0$ is the mass per unit length, $p(x)>0$ is the modulus of elasticity multiplied by the cross-sectional area (see [13] p.291), $\mu>0$ is a parameter, and the nonlinear forcing term $f(x, t, u)$ is $(2 \pi / \omega)$-periodic in time.

Equation (1) is a nonlinear model also for propagation of waves in nonisotropic media describing seismic phenomena, see e.g. [2].

We look for $(2 \pi / \omega)$-time periodic solutions $u(x, t)$ of (1).

[^0]This problem has received wide attention starting from the pioneering paper of Rabinowitz [21] dealing with the weakly nonlinear homogeneous string $\rho(x)=p(x)=1, \mu$ small, and the forcing frequency $\omega=1$ which enters in resonance with the proper eigen-frequencies $\omega_{j}=j \in \mathbb{N}$ of the string. For functions $2 \pi$-periodic in time and satisfying spatial Dirichlet boundary conditions, the spectrum $\left\{l^{2}-j^{2}, l \in \mathbb{Z}, j \geq 1\right\}$ of the D'Alembertian operator $\partial_{t t}-\partial_{x x}$ possesses the zero eigenvalue with infinite multiplicity (resonance) but the remaining eigenvalues are well separated. The corresponding infinite dimensional bifurcation problem is solved in [21] for nonlinearities $f$ which are monotone in $u$; see [6] and references therein for non-monotone $f$.

Subsequently many other results, both of bifurcation and of global nature (i.e. $\mu=1$ ), have been obtained, still for rational forcing frequencies $\omega \in \mathbb{Q}$, relying on these good separation properties of the spectrum, see e.g. [22, 23, $12,26,4]$ and references therein.

When the forcing frequency $\omega \in \mathbb{R} \backslash \mathbb{Q}$ is irrational (non-resonant case) the situation is completely different. Indeed the D'Alembertian operator $\omega^{2} \partial_{t t}-\partial_{x x}$ does not possess the zero eigenvalue but its spectrum $\left\{\omega^{2} l^{2}-j^{2}\right.$, $l \in \mathbb{Z}, j \geq 1\}$ accumulates to zero for almost every $\omega$. This is a "small divisors problem".

We underline that this "small divisors" phenomenon arises naturally for more realistic model equations like (1) where the density $\rho(x)$ and the modulus of elasticity $p(x)$ are not constant. Indeed in this case the eigenfrequencies $\omega_{j}$ of the string are no more integer numbers, having the asymptotic expansion

$$
\begin{equation*}
\omega_{j}^{2}=\frac{j^{2}}{c^{2}}+b+O\left(\frac{1}{j}\right) \tag{2}
\end{equation*}
$$

with suitable constants $c, b$ depending on $\rho, p$, see (55).
If $\omega=m / n \in \mathbb{Q}$, good separation properties of the spectrum can been recovered when $p(x)=\rho(x)$ (so $c=1$ ) and assuming the extra condition $b \neq 0$, see $[3,24]$. Indeed in this case the linear spectrum

$$
\omega^{2} l^{2}-\omega_{j}^{2}=\omega^{2} l^{2}-j^{2}-b+O\left(\frac{1}{j}\right)
$$

possesses at most finitely many zero eigenvalues and the remaining part of the spectrum is far away from zero. On the other hand, if $b=0$, the eigenvalues with $(l, j) \in(n, m) \mathbb{Z}$ tend to zero (also in the case $\omega \in \mathbb{Q}!$ ).

Existence of weak solutions in the non-resonant case was proved by Acquistapace [1] for $\rho=1$, for weak nonlinearities (i.e. $\mu$ small), and for a zero measure set of forcing frequencies $\omega$ for which the eigenvalues $\omega^{2} l^{2}-\omega_{j}^{2}$
are far away from zero. These frequencies are essentially the numbers whose continued fraction expansion is bounded, see [25].

For a similar zero measure set of frequencies, McKenna [18] has obtained some result when $\mu=1$, for $\rho=p=1$, and $f(t, x, u)=g(u)+h(t, x)$ with $g$ uniformly Lipschitz, via a fixed point argument, see also [5]; see [16, 9] for related results using variational methods.

Existence of classical solutions of (1) for a positive measure set of frequencies was proved by Plotnikov-Youngerman [19] for the homogeneous string $\rho=p=1, \mu$ small, and $f$ monotone in $u$. The monotonicity condition allows to control the first coefficient in the asymptotic expansion of the eigenvalues (as in (2)) of some perturbed linearized operator.

Recently Fokam [17] has proved existence of classical periodic solutions for large frequencies $\omega$ in a set of asymptotically full measure, for the homogeneous string $\rho=p=1$ plus a potential, when $\mu=1$ and $f=u^{3}+h(t, x)$ with $h$ a trigonometric polynomial odd in time and space.

In the present paper we prove existence of classical solutions of the nonhomogeneous string (1) for every $\rho(x), p(x)>0$, for general nonlinear terms $f(x, t, u)$, and for $(\mu, \omega)$ belonging to a large measure Cantor set $B_{\gamma}$, when the ratio $\mu / \omega$ is small, see Theorem 1, covering both the case $\mu \rightarrow 0$ and the case $\omega \rightarrow+\infty$.

In the limit $\mu / \omega \rightarrow 0$ the solution we find tends to a static equilibrium $v(x)$ with smaller, zero average oscillations $w(x, t)$ of amplitude $O(\mu / \omega)$, see (9),(10) and figure 1. The nonlinearity $f$ selects such $v$ through the infinite dimensional bifurcation equation (7) which possesses non-degenerate solutions under natural assumptions on $f$, see hypothesys (V). This problem was not present in [17] thanks to the symmetry assumptions on $f$.


Figure 1: The solution $u(x, t)=v(x)+w(x, t)$ of (1).
Considering the structure of the expected solution it is natural to attack the problem via a Lyapunov-Schmidt decomposition.

In the range equation (to find $w$ ) a small divisors problem arises and we solve it with a Nash-Moser type iterative scheme. The inversion of the
"linearized operators" - which is the core of any Nash-Moser scheme - is obtained adapting the technique of [7] to the present time-dependent case (section 4). See also $[10,11,14,15]$ for similar techniques.

It is here where the interaction between the forcing frequency $\omega$ and the normal modes of oscillations of the string linearized at different positions (approximating better and better the final string configuration) appears. The set $B_{\gamma}$ of "non-resonant" parameters $(\mu, \omega)$ for which we find a solution of the range equation (and then of (1)) is constructed avoiding these primary resonances. In particular the forcing frequency $\omega$ must not enter in resonance with the normal frequencies of oscillations of the string linearized at the limiting solution, see (11). At the end of the construction we obtain a large measure Cantor set $B_{\gamma}$ which looks like in figure 2. Outside this set the effect of resonance phenomena shall in general destroy the existence of periodic solutions like those found in Theorem 1.

We now present rigorously our results.

### 1.1 Main result

After a time rescaling we look for $2 \pi$-periodic solutions of

$$
\left\{\begin{array}{l}
\omega^{2} \rho(x) u_{t t}-\left(p(x) u_{x}\right)_{x}=\mu f(x, t, u)  \tag{3}\\
u(0, t)=u(\pi, t)=0
\end{array}\right.
$$

where $\mu \in[0, \bar{\mu}]$ for some $\bar{\mu}>0$, under the $2 \pi$-periodic forcing term

$$
\begin{equation*}
f(x, t, u)=\sum_{l \in \mathbb{Z}} f_{l}(x, u) e^{i l t}=f_{0}(x, u)+\bar{f}(x, t, u) \tag{4}
\end{equation*}
$$

where $\bar{f}(x, t, u):=\sum_{l \neq 0} f_{l}(x, u) e^{i l t}$.
We suppose that $f$ is analytic in $(t, u)$ :

$$
f(x, t, u)=\sum_{l \in \mathbb{Z}, k \in \mathbb{N}} f_{l k}(x) u^{k} e^{i l t}
$$

where $f_{l k}(x) \in H^{1}((0, \pi) ; \mathbb{C})$ and $^{2} f_{-l, k}=f_{l k}^{*}$.
Hypothesis (F). There exist $\sigma_{0}>0, r>0$ such that

$$
\sum_{l \in \mathbb{Z}}\left\|f_{l k}\right\|_{H^{1}}^{2}\left(1+l^{2}\right) e^{\left(2 \sigma_{0}\right) 2|l|}:=C_{k}^{2}(f)<\infty \quad \text { and } \quad \sum_{k=0}^{+\infty} C_{k}(f) r^{k}<\infty
$$

[^1]For example, any nonlinearity $f(x, t, u)$ which is a trigonometric polynomial in $t$ and a polynomial in $u$ satisfies hypothesys (F) for every $\sigma_{0}, r$.

If $f(x, t, 0) \neq 0$ equation (3) does not possess the trivial solution $u=0$.
We look for periodic solutions of (3) in the Hilbert space

$$
\begin{gathered}
X_{\sigma, s}:=\left\{u: \mathbb{T} \rightarrow H_{0}^{1}((0, \pi) ; \mathbb{R}), u(x, t)=\sum_{l \in \mathbb{Z}} u_{l}(x) e^{i l t}, u_{l} \in H_{0}^{1}((0, \pi) ; \mathbb{C}),\right. \\
\left.u_{-l}=u_{l}^{*},\|u\|_{\sigma, s}^{2}:=\sum_{l \in \mathbb{Z}}\left\|u_{l}\right\|_{H^{1}}^{2}\left(1+l^{2 s}\right) e^{2 \sigma|l|}<\infty\right\}
\end{gathered}
$$

of $2 \pi$-periodic in time functions valued in $H^{1}((0, \pi) ; \mathbb{R})$ which have a bounded analytic extension on the complex strip $|\operatorname{Im} t|<\sigma$ with trace function on $|\operatorname{Im} t|=\sigma$ belonging to $H^{s}\left(\mathbb{T} ; H^{1}((0, \pi) ; \mathbb{C})\right)$. For $s>1 / 2, X_{\sigma, s}$ is an algebra:

$$
\|u v\|_{\sigma, s} \leq c_{s}\|u\|_{\sigma, s}\|v\|_{\sigma, s} \quad \forall u, v \in X_{\sigma, s}
$$

with

$$
c_{s}:=2^{s}\left(\sum_{n \in \mathbb{Z}} \frac{1}{1+n^{2 s}}\right)^{1 / 2} .
$$

We shall use the notation $X_{\sigma}$, resp. $\left\|\|_{\sigma}\right.$, for $X_{\sigma, 1}$, resp. $\| \|_{\sigma, 1}$.
To find solutions of (3) we implement the Lyapunov-Schmidt reduction according to the decomposition

$$
X_{\sigma, s}=V \oplus\left(W \cap X_{\sigma, s}\right)
$$

where

$$
V:=H_{0}^{1}(0, \pi), \quad W:=\left\{w=\sum_{l \neq 0} w_{l}(x) e^{i l t} \in X_{0, s}\right\}
$$

writing every $u \in X_{\sigma, s}$ as $u=u_{0}(x)+\sum_{l \neq 0} u_{l}(x) e^{i l t}$.
Projecting equation (3), with $u(x, t)=v(x)+w(x, t), v \in V, w \in W$, yields

$$
\begin{cases}-\left(p v^{\prime}\right)^{\prime}=\mu \Pi_{V} f(v+w) & \text { bifurcation equation }  \tag{5}\\ L_{\omega} w=\mu \Pi_{W} f(v+w) & \text { range equation }\end{cases}
$$

where $\Pi_{V}, \Pi_{W}$ denote the projectors, $f(u)(x, t):=f(x, t, u(x, t))$ and

$$
L_{\omega} u:=\omega^{2} \rho(x) u_{t t}-\left(p(x) u_{x}\right)_{x}
$$

We shall find solutions of (5) when the ratio $\mu / \omega$ is small. In this limit $w$ tends to 0 and the bifurcation equation reduces to the time-independent equation

$$
\begin{equation*}
-\left(p v^{\prime}\right)^{\prime}=\mu f_{0}(v) \tag{6}
\end{equation*}
$$

because, by (4), for $w=0$,

$$
\Pi_{V} f(v)=\Pi_{V} f_{0}(x, v(x))+\Pi_{V} \bar{f}(x, t, v(x))=f_{0}(v) .
$$

The infinite dimensional " 0 -th order bifurcation equation" (6) is a nonautonomous second order ordinary differential equation, which, under natural conditions on $f_{0}$, possesses non-degenerate solutions satisfying the boundary conditions $v(0)=v(\pi)=0$.

Hypothesys (V). The problem

$$
\left\{\begin{array}{l}
-\left(p(x) v^{\prime}(x)\right)^{\prime}=\mu f_{0}(x, v(x))  \tag{7}\\
v(0)=v(\pi)=0
\end{array}\right.
$$

admits a solution $\bar{v} \in H_{0}^{1}(0, \pi)$ which is non-degenerate, namely the linearized equation

$$
-\left(p h^{\prime}\right)^{\prime}=\mu f_{0}^{\prime}(\bar{v}) h
$$

possesses in $H_{0}^{1}(0, \pi)$ only the trivial solution $h=0$.
We note that for $\mu=0$, the trivial solution $\bar{v}=0$ is non-degenerate, so, by the implicit function theorem, Hypothesis (V) is automatically satisfied for $\mu$ small. We deal also with the case $\mu$ not small, see for example Lemmas 2 and 3.

For the difficulties with a possibly degenerate solution we refer to [8].
Let $\lambda_{j}$ denote the eigenvalues of the Sturm-Liouville problem

$$
\left\{\begin{array}{l}
-\left(p(x) y^{\prime}(x)\right)^{\prime}=\lambda \rho(x) y(x)  \tag{8}\\
y(0)=y(\pi)=0
\end{array}\right.
$$

and $\omega_{j}:=\sqrt{\lambda_{j}}$. These are the frequencies of the free vibrations of the string (note that all the eigenvalues $\lambda_{j}$ are positive). Physically, it is the sequence of the fundamental tone $\omega_{1}$ and all its overharmonics $\omega_{2}, \omega_{3}, \ldots$ which compose the musical note of the string.

Theorem 1. Suppose $p(x), \rho(x)>0$ are of class $H^{3}(0, \pi), f$ satisfies $(F)$ and hypothesys ( $V$ ) holds for some $\mu_{0} \in[0, \bar{\mu}]$. Consider the open set
$A_{0}:=\left\{(\mu, \omega) \in\left(\mu_{1}, \mu_{2}\right) \times(\gamma,+\infty):\left|\omega l-\omega_{j}\right|>\frac{\gamma}{l^{\tau}}, \forall l=1, \ldots, N_{0}, j \geq 1\right\}$
where $\omega_{j}$ are defined by (8), $\gamma \in(0,1), \tau \in(1,2),\left(\mu_{1}, \mu_{2}\right)$ is a neighborhood of $\mu_{0}$ (see Lemma 4) and $N_{0}=N_{0}(\rho, p, f, \bar{\mu}, \bar{v}, \tau) \in \mathbb{N}$ is fixed in Lemma 7.
(Existence) There are constants $C, C^{\prime}>0$ depending only on $\rho, p, f, \bar{\mu}, \bar{v}, \tau$, a $\mathcal{C}^{\infty}$ function

$$
\tilde{w}: \tilde{A}:=A_{0} \cap\left\{(\mu, \omega): \frac{\mu}{\omega} \leq C^{\prime} \gamma^{5}\right\} \rightarrow X_{\sigma_{0} / 2} \cap W
$$

and a large, see section 3.3, Cantor set $B_{\gamma} \subset \tilde{A}$, such that for every $(\mu, \omega) \in$ $B_{\gamma}$ there exists a classical solution of (3)

$$
\begin{equation*}
\tilde{u}(\mu, \omega)=v(\mu, \tilde{w}(\mu, \omega))+\tilde{w}(\mu, \omega) \in V \oplus\left(W \cap X_{\sigma_{0} / 2}\right) \tag{9}
\end{equation*}
$$

satisfying

$$
\begin{equation*}
\|\tilde{w}(\mu, \omega)\|_{\sigma_{0} / 2} \leq C \frac{\mu}{\gamma \omega}, \quad\|v(\mu, \tilde{w}(\mu, \omega))-v(\mu, 0)\|_{H^{1}} \leq C \frac{\mu}{\gamma \omega} \tag{10}
\end{equation*}
$$

and $\|v(\mu, 0)-\bar{v}\|_{H^{1}} \leq C\left|\mu-\mu_{0}\right|$. The Cantor set $B_{\gamma}$ is explicitely

$$
\begin{align*}
B_{\gamma}:= & \left\{(\mu, \omega) \in\left(\mu_{1}, \mu_{2}\right) \times(2 \gamma,+\infty):\left|\omega l-\omega_{j}\right|>\frac{2 \gamma}{l^{\tau}} \quad \forall l=1, \ldots, N_{0}, j \geq 1,\right. \\
& \left.\frac{\mu}{\omega} \leq C^{\prime} \gamma^{5}, \quad\left|\omega l-\frac{j}{c}\right|>\frac{2 \gamma}{l^{\tau}}, \quad\left|\omega l-\tilde{\omega}_{j}(\mu, \omega)\right|>\frac{2 \gamma}{l^{\tau}} \quad \forall l, j \geq 1\right\} \quad(11) \tag{11}
\end{align*}
$$

where

$$
\begin{equation*}
c:=\frac{1}{\pi} \int_{0}^{\pi}\left(\frac{\rho(x)}{p(x)}\right)^{1 / 2} d x \tag{12}
\end{equation*}
$$

and $\tilde{\lambda}_{j}(\mu, \omega):=\tilde{\omega}_{j}^{2}(\mu, \omega)$ denote the (possibly negative ${ }^{3}$ ) eigenvalues of the Sturm-Liouville problem

$$
\left\{\begin{array}{l}
-\left(p y^{\prime}\right)^{\prime}-\mu \Pi_{V} f^{\prime}(v(\mu, \tilde{w}(\mu, \omega))+\tilde{w}(\mu, \omega)) y=\lambda \rho y  \tag{13}\\
y(0)=y(\pi)=0
\end{array}\right.
$$

(Regularity) Suppose, furthermore, $\rho(x) \in H^{m}(0, \pi), p(x) \in H^{m+1}(0, \pi)$, $f_{l k}(x) \in H^{m}(0, \pi)$ and $\sum_{l, k}\left\|f_{l k}\right\|_{H^{m}} r_{m}^{k}<\infty$ for some $m \geq 3, r_{m}>0$. If $\|\tilde{u}(\cdot, t)\|_{H_{0}^{1}}<r_{m} / C_{m}$ for some $C_{m}>0$, then $\tilde{u}(\cdot, t) \in H_{0}^{1}(0, \pi) \cap H^{m+2}(0, \pi)$.

This conclusion holds true, for example, when $f_{0}(x, 0)=d_{u} f_{0}(x, 0)=0$ (so $v(\mu, 0)=0, \forall \mu$ ) for $\mu / \gamma \omega$ small enough, by (10).

Fixed $\mu$, for every frequency $\omega$ in the section

$$
S(\mu):=\left\{\omega:(\mu, \omega) \in \cup_{\gamma \in(0,1)} B_{\gamma}\right\}
$$

[^2]

Figure 2: The Cantor set $B_{\gamma}$.
there exists a solution of (1) by Theorem 1. $S(\mu)$ has asymptotically full measure at $\omega \rightarrow+\infty$, i.e.

$$
\lim _{\omega \rightarrow+\infty}|S(\mu) \cap(\omega, \omega+1)|=1 .
$$

Analogously, fixed $\omega$, for every $\mu$ in the section

$$
S(\omega):=\left\{\mu:(\mu, \omega) \in \cup_{\gamma \in(0,1)} B_{\gamma}\right\}
$$

there exists a solution of (1). For $\mu$ small enough, also $S(\omega)$ is a "large" set: for every $\gamma^{\prime} \in(0,1), \omega^{\prime}>0$,

$$
\lim _{\mu \rightarrow 0}\left|\left\{\omega \in\left(\omega^{\prime}, \omega^{\prime}+1\right): \frac{|S(\omega) \cap(0, \mu)|}{\mu} \geq 1-\gamma^{\prime}\right\}\right|=1
$$

see section 3.3.
Notations. The symbols $K, K_{i}$ shall denote positive constants depending only on $\rho, p, f, \bar{\mu}, \bar{v}, \tau$.

## 2 The bifurcation equation

We first prove the analyticity of the Nemitski operator induced by $f$.

Lemma 1. Let $f$ satisfy assumption ( $F$ ). For every $\sigma \in\left[0, \sigma_{0}\right]$, $s>1 / 2$, the Nemitski operator $f$ is analytic on the ball $\left\{u \in X_{\sigma, s}: c_{s}\|u\|_{\sigma, s}<r\right\}$.

Proof. First note that
$\sum_{l \in \mathbb{Z}}\left\|u_{l}\right\|_{\infty} \leq \sqrt{\frac{\pi}{2}} \sum_{l \in \mathbb{Z}}\left\|u_{l}\right\|_{H^{1}} \leq \sqrt{\frac{\pi}{2}}\left(\sum_{l \in \mathbb{Z}}\left\|u_{l}\right\|_{H^{1}}^{2}\left(1+l^{2 s}\right)\right)^{1 / 2}\left(\sum_{l \in \mathbb{Z}} \frac{1}{1+l^{2 s}}\right)^{1 / 2}$
so $\|u\|_{\infty} \leq c_{s}\|u\|_{\sigma, s}, \forall u \in X_{\sigma, s}, \sigma \geq 0, s>1 / 2$, and $f(x, t, u(x, t))$ is welldefined.

By definition of the norm $\left\|\|_{\sigma, s}\right.$, there exists $C:=C\left(\sigma_{0}, s\right)>0$ such that $\forall \sigma \in\left[0, \sigma_{0}\right], \forall k \in \mathbb{N}$,

$$
\left\|\sum_{l \in \mathbb{Z}} f_{l k}(x) e^{i l t}\right\|_{\sigma, s} \leq C\left\|\sum_{l \in \mathbb{Z}} f_{l k}(x) e^{i l t}\right\|_{2 \sigma_{0}, 1}=C C_{k}(f)<+\infty
$$

by the assumption (F). Hence, for $c_{s}\|u\|_{\sigma, s}<r$, using the algebra property of $X_{\sigma, s}$,

$$
\begin{aligned}
\|f(u)\|_{\sigma, s} & \leq \sum_{k=0}^{\infty}\left\|\left(\sum_{l \in \mathbb{Z}} f_{l k}(x) e^{i l t}\right) u^{k}\right\|_{\sigma, s} \leq C \sum_{k=0}^{\infty} C_{k}(f)\left(c_{s}\|u\|_{\sigma, s}\right)^{k} \\
& <C \sum_{k=0}^{\infty} C_{k}(f) r^{k}<+\infty
\end{aligned}
$$

again by (F). The analyticity of the Nemitski operator $f$ w.r.t. $\left\|\|_{\sigma, s}\right.$ follows from the properties of the power series (see e.g. [20], Appendix A).

Throughout this paper we shall use the spaces $X_{\sigma, s}$ with $\sigma \in\left[\sigma_{0} / 2, \sigma_{0}\right]$ and $s \in \mathcal{S}:=\left\{1,1-\frac{\tau-1}{2}, 1+\frac{(\tau-1) \tau}{2-\tau}\right\}$. So we can choose a multiplicative algebra constant on $X_{\sigma, s}$ and a radius $R_{0}$ such that in the ball $\left\{u \in X_{\sigma, s}\right.$ : $\left.\|u\|_{\sigma, s}<R_{0}\right\} f$ is analytic, and $f, f^{\prime}, f^{\prime \prime}, \ldots$ are bounded, uniformly in $\sigma, s$.

We now give an example in which hypothesis $(V)$ holds.
Lemma 2. Suppose $f_{0}(x, u)=u^{m}$ for $m \geq 3$ odd and $p(x) \equiv 1$. Then, $\forall \mu$, there exists an unbounded sequence of non-degenerate solutions $v_{n}$ of (7).

Proof. All the solutions of the autonomous equation $-v^{\prime \prime}=\mu v^{m}$ are periodic and can be parametrized by their energy

$$
E=\frac{1}{2} v^{\prime 2}+\frac{\mu}{m+1} v^{m+1} .
$$

We denote $T_{E}$ the period of the solution $v_{E}$. We can suppose $v_{E}(0)=0$, so $v_{E}^{\prime}(0)=\sqrt{2 E}$. The other boundary condition $v_{E}(\pi)=0$ is satisfied iff

$$
\begin{equation*}
k \frac{T_{E}}{2}=\pi \quad \text { for some } k \in \mathbb{N} \tag{14}
\end{equation*}
$$

By symmetry and energy conservation $v_{E}\left(T_{E} / 4\right)=[(m+1) E / \mu]^{\frac{1}{m+1}}$. So

$$
\begin{aligned}
T_{E} & =4 \int_{0}^{\left[\frac{(m+1) E}{\mu}\right]^{\frac{1}{m+1}}}\left[2\left(E-\frac{\mu x^{m+1}}{m+1}\right)\right]^{-1 / 2} d x \\
& =\frac{4(m+1 / \mu)^{\frac{1}{m+1}}}{E^{\frac{1}{2}-\frac{1}{m+1}}} \int_{0}^{1} \frac{d y}{\sqrt{2\left(1-y^{m+1}\right)}}=\frac{C(m, \mu)}{E^{\frac{1}{2}-\frac{1}{m+1}}}
\end{aligned}
$$

by the change of variable $y=x[E(m+1) / \mu]^{-\frac{1}{m+1}}$, and (14) is satisfied at infinitely many energy levels. Let $\bar{E}>0$ such that $T_{\bar{E}}=2 \pi / k$ and denote the solution $\bar{v}:=v_{\bar{E}}$.

Let us prove that $\bar{v}$ is non-degenerate. Any solution $h$ of the linearized equation at $\bar{v}$,

$$
\begin{equation*}
-h^{\prime \prime}(x)=\mu m \bar{v}^{m-1}(x) h(x), \tag{15}
\end{equation*}
$$

can be written as $h=A \bar{v}^{\prime}+B \beta, A, B \in \mathbb{R}$, because $\bar{v}^{\prime}(x)$ and $\beta(x):=$ $\left(\partial_{E} v_{E}\right)_{\mid E=\bar{E}}(x)$ are solutions of (15); they are independent because $\bar{v}^{\prime}(0) \neq 0$ while $\beta(0)=0$. If $h(0)=0$ then $A=0$. We claim that $\beta(\pi) \neq 0$; as a consequence, if $h(\pi)=0$, then $B=0$, and so $h=0$, i.e. $\bar{v}$ is non-degenerate. To prove that $\beta(\pi) \neq 0$, we differentiate at $\bar{E}$ the identity $v_{E}\left(k T_{E} / 2\right)=0$,

$$
\beta(\pi)+\bar{v}^{\prime}(\pi)\left(\partial_{E} T_{E}\right)_{\mid E=\bar{E}}=0 .
$$

Since $\bar{v}^{\prime}(\pi)=(-1)^{k} \sqrt{2 E} \neq 0$ and $\partial_{E} T_{E} \neq 0$, we get $\beta(\pi) \neq 0$.

Lemma 3. If $f_{0}(x, 0)=d_{u} f_{0}(x, 0)=0$, then $\bar{v}=0$ is a non-degenerate solution of (7) for every $\mu$.

Proof. The linearized equation $-\left(p h^{\prime}\right)^{\prime}=0, h(0)=h(\pi)=0$ has only the trivial solution.

When hypothesys (V) holds at some $\left(\mu_{0}, \bar{v}\right)$, we solve first the bifurcation equation in (5) using the standard implicit function theorem. We find, for every $w$ small enough and $\mu$ in a neighborhood of $\mu_{0}$, a unique solution $v(\mu, w)$ of the bifurcation equation.

Lemma 4. There exist $0<R<R_{0}$, a neighborhood $\left[\mu_{1}, \mu_{2}\right.$ ] of $\mu_{0}$ and, $\forall \sigma \in\left[\sigma_{0} / 2, \sigma_{0}\right], s \in \mathcal{S}, a \mathcal{C}^{\infty} m a p$

$$
\left[\mu_{1}, \mu_{2}\right] \times\left\{w \in W \cap X_{\sigma, s}:\|w\|_{\sigma, s}<R\right\} \rightarrow V, \quad(\mu, w) \mapsto v(\mu, w)
$$

which solves the bifurcation equation in (5).
Proof. The linear operator

$$
h \mapsto-\left(p h^{\prime}\right)^{\prime}-\mu_{0} d_{v} \Pi_{V} f(v)[h]=-\left(p h^{\prime}\right)^{\prime}-\mu_{0} f_{0}^{\prime}(v) h
$$

is invertible on $H_{0}^{1}(0, \pi)$ by hypothesys $(\mathrm{V})$.

Remark 1. The solutions of the 0 th order bifurcation equation (7) found in Lemmas 2 and 3 are non-degenerate for every $\mu$, so, in such a case, we can continue $v(\mu, w)$ for all $\left[\mu_{1}, \mu_{2}\right]=[0, \bar{\mu}]$.

We denote by $\lambda_{j}(\mu, w):=\omega_{j}^{2}(\mu, w)$ the possibly negative eigenvalues of the Sturm-Liouville problem

$$
\left\{\begin{array}{l}
-\left(p y^{\prime}\right)^{\prime}-\mu \Pi_{V} f^{\prime}(v(\mu, w)+w) y=\lambda \rho y  \tag{16}\\
y(0)=y(\pi)=0
\end{array}\right.
$$

By a comparison principle, see the Appendix, the eigenvalues of (16) satisfy

$$
\begin{equation*}
\left|\lambda_{j}(\mu, w)-\lambda_{j}\left(\mu^{\prime}, w^{\prime}\right)\right| \leq K\left(\left|\mu-\mu^{\prime}\right|+\left\|w-w^{\prime}\right\|_{\sigma, s}\right) . \tag{17}
\end{equation*}
$$

The non-degeneracy of $\bar{v}:=v\left(\mu_{0}, 0\right)$ means that $\lambda_{j}\left(\mu_{0}, 0\right) \neq 0 \forall j$ and by (17)

$$
\begin{equation*}
\delta_{0}:=\inf \left\{\left|\lambda_{j}(\mu, w)\right|: j \geq 1, \quad \mu \in\left[\mu_{1}, \mu_{2}\right],\|w\|_{\sigma_{0} / 2} \leq R\right\}>0 \tag{18}
\end{equation*}
$$

eventually taking $R$ smaller. Note also that the index $j_{0}$ of the smallest positive eigenvalue is constant, independently on ( $\mu, w$ ).

## 3 Solution of the range equation

It remains to solve the range equation

$$
\begin{equation*}
L_{\omega} w=\mu \Pi_{W} \mathcal{F}(\mu, w) \tag{19}
\end{equation*}
$$

where

$$
\mathcal{F}(\mu, w):=f(v(\mu, w)+w) .
$$

By the previous lemmas, $\mathcal{F}$ is $\mathcal{C}^{\infty}$ and bounded, togheter with its derivatives, on $\left[\mu_{1}, \mu_{2}\right] \times B_{R}$ where $B_{R}:=\left\{w \in W \cap X_{\sigma, s}:\|w\|_{\sigma, s}<R\right\}$.

### 3.1 The Nash-Moser recursive scheme

We define the sequence of finite-dimensional subspaces

$$
W^{(n)}:=\left\{w=\sum_{1 \leq|l| \leq N_{n}} w_{l}(x) e^{i l t}\right\} \subset W
$$

where $N_{n}:=N_{0} 2^{n}$ and $N_{0} \in \mathbb{N}$. We also set

$$
W^{(n) \perp}:=\left\{w=\sum_{\left|| |>N_{n}\right.} w_{l}(x) e^{i l t} \in W\right\}
$$

and denote $P_{n}$ the projection on $W^{(n)}, P_{n}^{\perp}$ on $W^{(n) \perp}$. For $w \in W^{(n) \perp}$ the following smoothing estimate holds: if $0<\sigma^{\prime \prime}<\sigma^{\prime}$,

$$
\begin{equation*}
\|w\|_{\sigma^{\prime \prime}, s} \leq \exp \left[-\left(\sigma^{\prime}-\sigma^{\prime \prime}\right) N_{n}\right]\|w\|_{\sigma^{\prime}, s} . \tag{20}
\end{equation*}
$$

The key property for the construction of the iterative sequence is the invertibility of the linear operator

$$
\begin{align*}
\mathcal{L}_{n}(w) h & :=-L_{\omega} h+\mu P_{n}\left[d_{w} \mathcal{F}(\mu, w) h\right]  \tag{21}\\
& =-L_{\omega} h+\mu P_{n}\left[f^{\prime}(v(\mu, w)+w)\left(h+d_{w} v(\mu, w)[h]\right)\right] \quad \forall h \in W^{(n)} .
\end{align*}
$$

Lemma 5. (Inversion of the linear problem) Let $\tau \in(1,2), \gamma \in(0,1)$, $\sigma \in\left(0, \sigma_{0}\right]$. Assume ${ }^{4} \omega>\gamma$ and the non-resonance conditions:

$$
\begin{equation*}
\left|\omega l-\frac{j}{c}\right|>\frac{\gamma}{l^{\tau}} \quad \forall l=1,2, \ldots, N_{n}, \quad \forall j \geq 1 \tag{22}
\end{equation*}
$$

where $c$ is defined in (12), and

$$
\begin{equation*}
\left|\omega^{2} l^{2}-\lambda_{j}(\mu, w)\right|>\frac{\gamma \omega}{l^{\tau-1}} \quad \forall l=1,2, \ldots, N_{n}, j \geq 1 \tag{23}
\end{equation*}
$$

where $\lambda_{j}(\mu, w)$ are the eigenvalues of (16).
Let $u:=v(\mu, w)+w$. There exist $K_{1}, K_{1}^{\prime}$ such that, if

$$
\begin{equation*}
\frac{\mu}{\gamma^{3} \omega}\left\|\Pi_{W} f^{\prime}(u)\right\|_{\sigma, 1+\frac{\tau(\tau-1)}{2-\tau}}<K_{1}^{\prime} \tag{24}
\end{equation*}
$$

then $\mathcal{L}_{n}(w)$ is invertible and

$$
\begin{equation*}
\left\|\mathcal{L}_{n}(w)^{-1} h\right\|_{\sigma} \leq \frac{K_{1} N_{n}^{\tau-1}}{\gamma \omega}\|h\|_{\sigma} \quad \forall h \in W^{(n)} \tag{25}
\end{equation*}
$$

[^3]Proof. In section 4.
For $\vartheta:=3 \sigma_{0} / \pi^{2}$ we define the sequence

$$
\begin{equation*}
\sigma_{n+1}:=\sigma_{n}-\frac{\vartheta}{(n+1)^{2}}, \quad \sigma_{0}>\sigma_{1}>\sigma_{2}>\ldots>\frac{\sigma_{0}}{2} \tag{26}
\end{equation*}
$$

Lemma 6. (The approximate solution) If $(\mu, \omega) \in A_{0}$ and $\mu N_{0}^{\tau-1} / \gamma \omega<$ $K_{2}^{\prime}$ is sufficiently small, then there exists a solution $w_{0}:=w_{0}(\mu, \omega) \in W^{(0)}$ of

$$
L_{\omega} w_{0}=\mu P_{0} \mathcal{F}\left(\mu, w_{0}\right)
$$

satisfying $\left\|w_{0}\right\|_{\sigma_{0}} \leq \mu K_{2} N_{0}^{\tau-1} / \gamma \omega$ for some $K_{2}$.
Proof. By definition of $A_{0}$ in Theorem 1, the eigenvalues of $(1 / \rho) L_{\omega}$ satisfy

$$
\left|\omega^{2} l^{2}-\lambda_{j}\right|>\frac{\gamma \omega}{l^{\tau-1}} \quad \forall l=1,2, \ldots, N_{0}, \quad \forall j \geq 1
$$

so $L_{\omega}$ is invertible on $W^{(0)}$ and, for some $K$,

$$
\begin{equation*}
\left\|L_{\omega}^{-1} h\right\|_{\sigma_{0}} \leq \frac{K N_{0}^{\tau-1}}{\gamma \omega}\|h\|_{\sigma_{0}} \quad \forall h \in W^{(0)} . \tag{27}
\end{equation*}
$$

Then we look for a solution $w_{0} \in W^{(0)}$ of $w_{0}=\mu L_{\omega}^{-1} P_{0} \mathcal{F}\left(\mu, w_{0}\right)$. The righthand side term is a contraction in $\left\{\left\|w_{0}\right\|_{\sigma_{0}}<R\right\}$ if $\mu N_{0}^{\tau-1} / \gamma \omega$ is sufficiently small.

Given $w_{n} \in W^{(n)},\left\|w_{n}\right\|_{\sigma_{n}}<R$ and $A_{n} \subseteq A_{0}$, we define

$$
\begin{aligned}
A_{n+1}:=\left\{(\mu, \omega) \in A_{n}:\left|\omega l-\omega_{j}\left(\mu, w_{n}\right)\right|\right. & >\frac{\gamma}{l^{\tau}}, \quad\left|\omega l-\frac{j}{c}\right|>\frac{\gamma}{l^{\tau}} \\
\forall l & \left.=1,2, \ldots, N_{n+1}, j \geq 1\right\} \subseteq A_{n}
\end{aligned}
$$

where $\lambda_{j}\left(\mu, w_{n}\right)=\omega_{j}^{2}\left(\mu, w_{n}\right)$ are defined in (16) with $w=w_{n}$.
In Lemma 6 we have constructed $h_{0}:=w_{0}$ for $(\mu, \omega) \in A_{0}$. Next, we proceed by induction. By means of $w_{0}$ we define the set $A_{1}$ as above, and we find $w_{1}:=h_{0}+h_{1} \in W^{(1)}$ for every $(\mu, \omega) \in A_{1}$ by Lemma 7 below. Then we define $A_{2}$, we find $w_{2} \in W^{(2)}$ and so on. The main goal of the construction is to prove that, at the end of the recurrence, the set of parameters $(\mu, \omega) \in$ $\cap_{n} A_{n}$ is actually a large set.
Lemma 7. (Inductive step). Fix $\chi \in(1,2)$. Suppose that $h_{i} \in W^{(i)}$, $\forall i=0, \ldots, n$, satisfy

$$
\begin{equation*}
\left\|h_{i}\right\|_{\sigma_{i}}<\frac{\mu K_{2} N_{0}^{\tau-1}}{\gamma \omega} \exp \left(-\chi^{i}\right) \tag{28}
\end{equation*}
$$

where $K_{2}$ is the constant in Lemma 6; $\forall k=0, \ldots, n, w_{k}:=h_{0}+\ldots+h_{k}$ satisfies $\left\|w_{k}\right\|_{\sigma_{k}}<R$ and

$$
\begin{equation*}
L_{\omega} w_{k}=\mu P_{k} \mathcal{F}\left(\mu, w_{k}\right) \tag{29}
\end{equation*}
$$

and suppose that $(\mu, \omega) \in A_{n}$, where $A_{i+1}$ is constructed by means of $w_{i}$ as showed above.

There exist $N_{0}=N_{0}(\rho, p, f, \bar{\mu}, \bar{v}, \tau) \in \mathbb{N}$ and $K_{3}^{\prime}$ such that: if $(\mu, \omega) \in$ $A_{n+1}$ and $\mu / \gamma^{3} \omega<K_{3}^{\prime}$, then there exists $h_{n+1} \in W^{(n+1)}$ satisfying

$$
\begin{equation*}
\left\|h_{n+1}\right\|_{\sigma_{n+1}}<\frac{\mu K_{2} N_{0}^{\tau-1}}{\gamma \omega} \exp \left(-\chi^{n+1}\right) \tag{30}
\end{equation*}
$$

such that $w_{n+1}=w_{n}+h_{n+1}$ verifies $\left\|w_{n+1}\right\|_{\sigma_{n+1}}<R$ and

$$
\begin{equation*}
L_{\omega} w_{n+1}=\mu P_{n+1} \mathcal{F}\left(\mu, w_{n+1}\right) \tag{31}
\end{equation*}
$$

Proof. In short $\mathcal{F}(w):=\mathcal{F}(\mu, w)$ and $D \mathcal{F}(w):=d_{w} \mathcal{F}(\mu, w)$. Equation (31) for $w_{n+1}=w_{n}+h_{n+1}$ is $L_{\omega}\left[w_{n}+h_{n+1}\right]=\mu P_{n+1} \mathcal{F}\left(w_{n}+h_{n+1}\right)$.

By assumption, $w_{n}$ satisfies (29) for $k=n$, namely $L_{\omega} w_{n}=\mu P_{n} \mathcal{F}\left(w_{n}\right)$, so the equation for $h_{n+1}$ can be written as

$$
\begin{equation*}
\mathcal{L}_{n+1}\left(w_{n}\right) h_{n+1}+\mu\left(P_{n+1}-P_{n}\right) \mathcal{F}\left(w_{n}\right)+\mu P_{n+1} Q=0 \tag{32}
\end{equation*}
$$

where, as defined in (21), $\mathcal{L}_{n+1}\left(w_{n}\right) h_{n+1}:=-L_{\omega} h_{n+1}+\mu P_{n+1} D \mathcal{F}\left(w_{n}\right) h_{n+1}$, and $Q$ denotes the quadratic remainder

$$
Q=Q\left(w_{n}, h_{n+1}\right):=\mathcal{F}\left(w_{n+1}\right)-\mathcal{F}\left(w_{n}\right)-D \mathcal{F}\left(w_{n}\right) h_{n+1} .
$$

Step 1: Inversion of $\mathcal{L}_{n+1}\left(w_{n}\right)$. We verify the assumptions of Lemma 5. By definition of $A_{n+1}, \omega$ satisfies (22). If $\lambda_{j}\left(\mu, w_{n}\right)<0$, then $\left|\omega^{2} l^{2}-\lambda_{j}\left(\mu, w_{n}\right)\right| \geq$ $\omega^{2} l^{2}>\gamma \omega / l^{\tau-1}$ because $\omega>\gamma$. If $\lambda_{j}\left(\mu, w_{n}\right)>0$, we have

$$
\left|\omega^{2} l^{2}-\lambda_{j}\left(\mu, w_{n}\right)\right| \geq\left|\omega l-\omega_{j}\left(\mu, w_{n}\right)\right| \omega l>\frac{\gamma \omega}{l^{\tau-1}} \quad \forall l=1, \ldots, N_{n+1}
$$

because $(\mu, \omega) \in A_{n+1}$. In both cases the non-resonance condition (23) holds.
To verify (24) we need an estimate for $w_{n}$. Let $\eta:=\tau(\tau-1) /(2-\tau)$ and $\alpha>0$. Using the elementary inequality

$$
\frac{1+l^{2(1+\eta)}}{1+l^{2}} \cdot \frac{e^{2(\sigma-\alpha)|l|}}{e^{2 \sigma|l|}} \leq \frac{2 l^{2 \eta}}{e^{2 \alpha|l|}} \leq 2\left(\frac{\eta}{\alpha e}\right)^{2 \eta}, \quad \forall l \neq 0
$$

we deduce

$$
\left\|h_{i}\right\|_{\sigma_{n+1}, 1+\eta} \leq \frac{C_{\eta}}{\left(\sigma_{i}-\sigma_{n+1}\right)^{\eta}}\left\|h_{i}\right\|_{\sigma_{i}}
$$

where $C_{\eta}:=\sqrt{2}(\eta / e)^{\eta}$. Since $\sigma_{i}-\sigma_{n+1} \geq \sigma_{i}-\sigma_{i+1}$ for every $i \leq n$,

$$
\left\|w_{n}\right\|_{\sigma_{n+1}, 1+\eta} \leq \sum_{i=0}^{n}\left\|h_{i}\right\|_{\sigma_{n+1}, 1+\eta} \leq C_{\eta} \sum_{i=0}^{n} \frac{\left\|h_{i}\right\|_{\sigma_{i}}}{\left(\sigma_{i}-\sigma_{i+1}\right)^{\eta}} \leq S_{\eta} \frac{\mu K_{2} N_{0}^{\tau-1}}{\gamma \omega}
$$

using (28) where $S_{\eta}:=\left(C_{\eta} / \vartheta^{\eta}\right) \sum_{i=0}^{+\infty}(i+1)^{2 \eta} \exp \left(-\chi^{i}\right)<+\infty$. If $S_{\eta} \mu K_{2}$ $N_{0}^{\tau-1} / \gamma \omega<R$ then $\left\|f^{\prime}\left(u_{n}\right)\right\|_{\sigma_{n+1}, 1+\eta} \leq K$ where $u_{n}:=v\left(\mu, w_{n}\right)+w_{n}$, and hypotheses (24) is verified for $\mu / \gamma^{3} \omega<K^{\prime}$ sufficiently small.

Analogously we get $\left\|w_{n}\right\|_{\sigma_{n}}<R$ if $\mu N_{0}^{\tau-1} / \gamma \omega<K^{\prime \prime}$ is small enough.
By Lemma 5 the operator $\mathcal{L}_{n+1}\left(w_{n}\right)$ is invertible on $W^{(n+1)}$ and

$$
\begin{equation*}
\left\|\mathcal{L}_{n+1}\left(w_{n}\right)^{-1} h\right\|_{\sigma_{n+1}} \leq \frac{K_{1} N_{n+1}^{\tau-1}}{\gamma \omega}\|h\|_{\sigma_{n+1}}, \quad \forall h \in W^{(n+1)} \tag{33}
\end{equation*}
$$

Equation (32) amounts to the fixed point problem

$$
h_{n+1}=-\mu \mathcal{L}_{n+1}\left(w_{n}\right)^{-1}\left[\left(P_{n+1}-P_{n}\right) \mathcal{F}\left(w_{n}\right)+P_{n+1} Q\right]:=\mathcal{G}\left(h_{n+1}\right)
$$

for $h_{n+1} \in W^{(n+1)}$.
Step 2: $\mathcal{G}$ is a contraction. We prove that $\mathcal{G}$ is a contraction on the ball $B_{n+1}:=\left\{\|h\|_{\sigma_{n+1}}<r_{n+1}\right\}$ where $r_{n+1}:=\left(\mu K_{2} N_{0}^{\tau-1} / \gamma \omega\right) \exp \left(-\chi^{n+1}\right)$, implying (30). By (20)

$$
\left\|\left(P_{n+1}-P_{n}\right) \mathcal{F}\left(w_{n}\right)\right\|_{\sigma_{n+1}} \leq\left\|\mathcal{F}\left(w_{n}\right)\right\|_{\sigma_{n}} \exp \left[-\left(\sigma_{n}-\sigma_{n+1}\right) N_{n}\right]
$$

Since $\left\|w_{n}\right\|_{\sigma_{n}}<R$, we have $\|Q\|_{\sigma_{n+1}} \leq K\left\|h_{n+1}\right\|_{\sigma_{n+1}}^{2}$. Hence, by (33),

$$
\left\|\mathcal{G}\left(h_{n+1}\right)\right\|_{\sigma_{n+1}} \leq K \frac{\mu N_{n+1}^{\tau-1}}{\gamma \omega}\left(\exp \left[-\left(\sigma_{n}-\sigma_{n+1}\right) N_{n}\right]+\left\|h_{n+1}\right\|_{\sigma_{n+1}}^{2}\right)
$$

Therefore $\mathcal{G}\left(B_{n+1}\right) \subseteq B_{n+1}$ if

$$
\begin{equation*}
\frac{\mu K N_{n+1}^{\tau-1}}{\gamma \omega} \exp \left[-\left(\sigma_{n}-\sigma_{n+1}\right) N_{n}\right]<\frac{r_{n+1}}{2}, \quad \frac{\mu K N_{n+1}^{\tau-1}}{\gamma \omega} r_{n+1}^{2}<\frac{r_{n+1}}{2} \tag{34}
\end{equation*}
$$

By the definition of $\sigma_{n}$ in (26) and $N_{n}:=N_{0} 2^{n}$, the first inequality is verified for every $n \geq 0$ if $\sigma_{0} N_{0}$ is greater than a constant depending only on $\chi, K, K_{2}$. The second inequality is verified for every $n \geq 0$ if $\mu N_{0}^{\tau-1} / \gamma \omega<K^{\prime}$ is small enough.

The estimate for $\|\mathcal{G} h-\mathcal{G} k\|, h, k \in B_{n+1}$ is similar. By the Contraction Mapping Theorem we conclude.

Corollary 1. (Existence) Suppose $A_{\infty}:=\bigcap_{n \geq 0} A_{n} \neq \emptyset$. If $(\mu, \omega) \in A_{\infty}$ then

$$
w_{\infty}(\mu, \omega):=\sum_{n \geq 0} h_{n}(\mu, \omega) \in W \cap X_{\sigma_{0} / 2}
$$

is a solution of the range equation (19) satisfying $\left\|w_{\infty}\right\|_{\sigma_{0} / 2} \leq K_{3} \mu / \gamma \omega$, and

$$
u_{\infty}:=v\left(\mu, w_{\infty}(\mu, \omega)\right)+w_{\infty}(\mu, \omega) \in X_{\sigma_{0} / 2}
$$

is a classical solution of (3) satisfying (10).
Proof. Since $w_{n}$ solves (29) for $k=n,-L_{\omega} w_{n}+\mu \Pi_{W} f\left(u_{n}\right)=\mu P_{n}^{\perp} f\left(u_{n}\right)$ $\in W^{(n) \perp}$ where $u_{n}:=v\left(\mu, w_{n}\right)+w_{n}$. By (20)

$$
\lim _{n \rightarrow+\infty}\left\|-L_{\omega} w_{n}+\mu f\left(u_{n}\right)\right\|_{\sigma_{0} / 2} \leq \lim _{n \rightarrow+\infty} K \exp \left[-\left(\sigma_{n}-\sigma_{0} / 2\right) N_{n}\right]=0
$$

Since $w_{n} \rightarrow w_{\infty}$ in $\left\|\|_{\sigma_{0} / 2}\right.$ also $f\left(u_{n}\right) \rightarrow f\left(u_{\infty}\right)$ in the same norm, while $L_{\omega} w_{n} \rightarrow L_{\omega} w_{\infty}$ in the sense of distributions. So $w_{\infty}$ is a weak solution of the range equation (19) and $u_{\infty}:=v\left(\mu, w_{\infty}(\mu, \omega)\right)+w_{\infty}(\mu, \omega) \in X_{\sigma_{0} / 2}$ is a weak solution of (3).

Finally, by the equation, $\partial_{x}\left(p(x) \partial_{x} u_{\infty}(x, t)\right)$ is a continuous function in $(x, t)$ and, $\forall t, u_{\infty}(\cdot, t) \in H^{3}(0, \pi) \subset \mathcal{C}^{2}$ is a classical solution of (3).

Remark 2. We shall prove, as a consequence of Lemma 11 and section 3.3, that $A_{\infty}$ is actually a positive measure set. A possible way to prove it uses the Whitney extension of $w_{\infty}$, see section 3.2.

Lemma 8. (Regularity) Suppose $\rho(x) \in H^{m}(0, \pi), p(x) \in H^{m+1}(0, \pi)$, $f_{l k}(x) \in H^{m}(0, \pi)$ and $\sum_{l, k}\left\|f_{l k}\right\|_{H^{m}} r_{m}^{k}<\infty$ for some $m \geq 3, r_{m}>0$.

There exists $K_{m}$ such that if the solution $u_{\infty}$ of Corollary 1 satisfies $\left\|u_{\infty}(\cdot, t)\right\|_{H^{1}}<K_{m}$, then $u_{\infty}(\cdot, t) \in H^{m+2}(0, \pi) \cap H_{0}^{1}(0, \pi)$.

Proof. For every fixed $t$, by the algebra property of $H^{m}$

$$
\|f(x, t, u(x, t))\|_{H^{m}} \leq \sum_{l, k}\left\|f_{l k}(x) u^{k}(x)\right\|_{H^{m}} \leq C \sum_{l, k}\left\|f_{l k}\right\|_{H^{m}}\left\|u^{k}\right\|_{H^{m}}
$$

Using the Gagliardo-Nirenberg type inequality $\left\|u^{k}\right\|_{H^{m}} \leq\left(C_{m}\|u\|_{H^{1}}\right)^{k-1}\|u\|_{H^{m}}$ valid for every $u \in H_{0}^{1} \cap H^{m}$, we get

$$
\begin{equation*}
\|f(x, t, u(x, t))\|_{H^{m}} \leq C\|u\|_{H^{m}} \sum_{l, k}\left\|f_{l k}\right\|_{H^{m}}\left(C_{m}\|u\|_{H^{1}}\right)^{k-1} \tag{35}
\end{equation*}
$$

which is convergent for $\|u\|_{H^{1}}<r_{m} / C_{m}$. The solution $u:=u_{\infty}$ satisfies

$$
\begin{equation*}
-\left(p(x) u_{x}\right)_{x}=\mu f(x, t, u)-\rho(x) u_{t t} \tag{36}
\end{equation*}
$$

and $u(\cdot, t) \in H^{3}(0, \pi), \forall t$. Suppose $\|u\|_{H^{1}}<r_{m} / C_{m}$. By induction, assume $u(\cdot, t) \in H^{k}$ for $k=3, \ldots, m$. Hence $u_{t t}(\cdot, t) \in H^{k}$ and by (35) $f(x, t, u) \in$ $H^{k}$. This implies by (36) that $u \in H^{k+2}$.

Remark 3. The solution $u_{\infty}$ is small if $v(\mu, 0)=0$, because $\left\|u_{\infty}\right\|_{H^{1}} \leq$ $\left\|u_{\infty}\right\|_{\sigma_{0} / 2} \leq K \mu / \gamma \omega$. In this case $u_{\infty}(\cdot, t) \in H^{m+2}$ for $\mu / \gamma \omega$ small enough.

### 3.2 Whitney $\mathcal{C}^{\infty}$ extension

The functions $h_{n}$ constructed in Lemmas 6 and 7 depend smoothly on the parameters $(\mu, \omega)$.

Lemma 9. There is $K_{4}, K_{4}^{\prime}$ such that for $\mu / \gamma^{3} \omega<K_{4}^{\prime}$, the map $h_{i}(\mu, \omega) \in$ $\mathcal{C}^{\infty}\left(A_{i}, W^{(i)}\right)$, and

$$
\left\|\partial_{\omega} h_{i}(\mu, \omega)\right\|_{\sigma_{i}} \leq \frac{K_{4} \mu}{\gamma^{2} \omega} \exp \left(-\chi_{0}^{i}\right), \quad\left\|\partial_{\mu} h_{i}(\mu, \omega)\right\|_{\sigma_{i}} \leq \frac{K_{4}}{\gamma \omega} \exp \left(-\chi_{0}^{i}\right)
$$

where $\chi_{0}:=(1+\chi) / 2$.
Proof. Since $w_{0}=\mu L_{\omega}^{-1} P_{0} \mathcal{F}\left(\mu, w_{0}\right)$, by the implicit function theorem the map $w_{0}(\mu, \omega) \in \mathcal{C}^{\infty}\left(A_{0}, W^{(0)}\right)$. Differentiating the identity $L_{\omega}\left(L_{\omega}^{-1} h\right)=h$ w.r.t. $\omega$, by (27), we get $\left\|\partial_{\omega} L_{\omega}^{-1} h\right\|_{\sigma_{0}} \leq\left(K / \gamma^{2} \omega\right)\|h\|_{\sigma_{0}}$. For $\mu / \gamma \omega$ small,

$$
\left\|\partial_{\omega} w_{0}\right\|_{\sigma_{0}} \leq \frac{K \mu}{\gamma^{2} \omega}
$$

Differentiating w.r.t. $\mu$ we get also $\left\|\partial_{\mu} w_{0}\right\|_{\sigma_{0}} \leq K^{\prime} / \gamma \omega$.
By induction, suppose that $h_{i}$ depends smoothly on $(\mu, \omega) \in A_{i}$ for every $i=0, \ldots, n$. For $(\mu, \omega) \in A_{n+1}$, by (31), $h_{n+1}$ is a solution of

$$
\begin{equation*}
-L_{\omega} h_{n+1}+\mu P_{n+1}\left[\mathcal{F}\left(w_{n}+h_{n+1}\right)-\mathcal{F}\left(w_{n}\right)\right]+\mu\left(P_{n+1}-P_{n}\right) \mathcal{F}\left(w_{n}\right)=0 \tag{37}
\end{equation*}
$$

By the implitic function theorem $h_{n+1} \in \mathcal{C}^{\infty}$ once we prove that

$$
\mathcal{L}_{n+1}\left(w_{n+1}\right)[z]:=-L_{\omega} z+\mu P_{n+1} D \mathcal{F}\left(w_{n}+h_{n+1}\right)[z]
$$

is invertible. By (33), $\mathcal{L}_{n+1}\left(w_{n}\right)$ is invertible. Hence it is sufficient that

$$
\left\|\mathcal{L}_{n+1}^{-1}\left(w_{n}\right)\left(\mathcal{L}_{n+1}\left(w_{n+1}\right)-\mathcal{L}_{n+1}\left(w_{n}\right)\right)\right\|_{\sigma_{n+1}}<\frac{1}{2},
$$

which holds true for $\mu^{2} / \gamma \omega$ small enough, since, by (30),

$$
\left\|\mathcal{L}_{n+1}\left(w_{n+1}\right)-\mathcal{L}_{n+1}\left(w_{n}\right)\right\|_{\sigma_{n+1}} \leq K \mu\left\|h_{n+1}\right\|_{\sigma_{n+1}} \leq \frac{\mu^{2} K^{\prime} N_{0}^{\tau-1}}{\gamma \omega} \exp \left(-\chi^{n+1}\right)
$$

Finally (33) implies

$$
\begin{equation*}
\left\|\mathcal{L}_{n+1}\left(w_{n+1}\right)^{-1}\right\|_{\sigma_{n+1}} \leq \frac{2 K_{1} N_{n+1}^{\tau-1}}{\gamma \omega} \tag{38}
\end{equation*}
$$

Differentiating (37) w.r.t. $\omega$

$$
\begin{array}{r}
\mathcal{L}_{n+1}\left(w_{n+1}\right)\left[\partial_{\omega} h_{n+1}\right]=2 \omega \rho(x)\left(h_{n+1}\right)_{t t}-\mu\left(P_{n+1}-P_{n}\right) D \mathcal{F}\left(w_{n}\right) \partial_{\omega} w_{n} \\
-\mu P_{n+1}\left[D \mathcal{F}\left(w_{n}+h_{n+1}\right)-D \mathcal{F}\left(w_{n}\right)\right] \partial_{\omega} w_{n} \tag{39}
\end{array}
$$

and, using (38) and (20),

$$
\begin{aligned}
\left\|\partial_{\omega} h_{n+1}\right\|_{\sigma_{n+1}} \leq & \frac{K N_{n+1}^{\tau-1}}{\gamma \omega}\left(\omega N_{n+1}^{2}\left\|h_{n+1}\right\|_{\sigma_{n+1}}+\frac{\mu\left\|\partial_{\omega} w_{n}\right\|_{\sigma_{n}}}{\exp \left[\left(\sigma_{n}-\sigma_{n+1}\right) N_{n}\right]}+\right. \\
& \left.+\mu\left\|h_{n+1}\right\|_{\sigma_{n+1}}\left\|\partial_{\omega} w_{n}\right\|_{\sigma_{n}}\right) .
\end{aligned}
$$

We note that $\left\|\partial_{\omega} w_{n}\right\|_{\sigma_{n}} \leq \sum_{i=0}^{n}\left\|\partial_{\omega} h_{i}\right\|_{\sigma_{i}}$. Using (34) the sequence $a_{n}:=$ $\left\|\partial_{\omega} h_{n}\right\|_{\sigma_{n}}$ satisfies

$$
\begin{aligned}
a_{n+1} & \leq \frac{K N_{n+1}^{\tau-1}}{\gamma \omega}\left(\omega N_{n+1}^{2} r_{n+1}+\frac{\omega \gamma r_{n+1}}{N_{n+1}^{\tau-1}} \sum_{i=0}^{n} a_{i}+\mu r_{n+1} \sum_{i=0}^{n} a_{i}\right) \\
& \leq b_{n+1}\left(1+\sum_{i=0}^{n} a_{i}\right) \text { where } b_{n+1}:=\frac{K \mu}{\gamma^{2} \omega} N_{n+1}^{\tau+1} \exp \left(-\chi^{n+1}\right),
\end{aligned}
$$

recalling that $r_{n+1}=(\mu K / \gamma \omega) \exp \left(-\chi^{n+1}\right)$. By induction, for $K \mu / \omega \gamma^{2}<1$, we have $a_{n} \leq 2 b_{n}$ and

$$
\left\|\partial_{\omega} h_{n+1}\right\|_{\sigma_{n+1}} \leq \frac{K \mu}{\gamma^{2} \omega} N_{n+1}^{\tau+1} \exp \left(-\chi^{n+1}\right) \leq \frac{K^{\prime} \mu}{\gamma^{2} \omega} \exp \left(-\chi_{0}^{n+1}\right)
$$

where $\chi_{0}:=(1+\chi) / 2$. It follows that $\left\|\partial_{\omega} w_{n+1}\right\|_{\sigma_{n+1}} \leq K \mu / \gamma^{2} \omega$.
Differentiating (37) w.r.t. $\mu$ we obtain the estimate for $\partial_{\mu} h_{n+1}$.
Define, for $\nu_{0}>0$,

$$
\begin{aligned}
& A_{n}^{*}:=\left\{(\mu, \omega) \in A_{n}: \operatorname{dist}\left((\mu, \omega), \partial A_{n}\right)>\frac{\nu_{0} \gamma^{4}}{N_{n}^{3}}\right\} \\
& \tilde{A}_{n}:=\left\{(\mu, \omega) \in A_{n}: \operatorname{dist}\left((\mu, \omega), \partial A_{n}\right)>\frac{2 \nu_{0} \gamma^{4}}{N_{n}^{3}}\right\} \subset A_{n}^{*}
\end{aligned}
$$

Lemma 10. (Whitney extension) There exists $\tilde{w} \in \mathcal{C}^{\infty}\left(A_{0}, W \cap X_{\sigma_{0} / 2}\right)$ satisfying

$$
\begin{equation*}
\|\tilde{w}(\mu, \omega)\|_{\sigma_{0} / 2} \leq \frac{K_{3} \mu}{\gamma \omega}, \quad\left\|\partial_{\omega} \tilde{w}(\mu, \omega)\right\|_{\sigma_{0} / 2} \leq \frac{C \mu}{\gamma^{5} \omega}, \quad\left\|\partial_{\mu} \tilde{w}(\mu, \omega)\right\|_{\sigma_{0} / 2} \leq \frac{C}{\gamma^{5} \omega} \tag{40}
\end{equation*}
$$

for some $C:=C\left(\nu_{0}\right)>0$, such that, $\forall(\mu, \omega) \in \tilde{A}_{\infty}:=\bigcap_{n \geq 0} \tilde{A}_{n}, \tilde{w}(\mu, \omega)$ solves the range equation (19).

Moreover there exists a sequence $\tilde{w}_{n} \in \mathcal{C}^{\infty}\left(A_{0}, W^{(n)}\right)$ such that $\tilde{w}_{n}(\mu, \omega)=$ $w_{n}(\mu, \omega), \forall(\mu, \omega) \in \tilde{A}_{n}$, and

$$
\begin{equation*}
\left\|\tilde{w}(\mu, \omega)-\tilde{w}_{n}(\mu, \omega)\right\|_{\sigma_{0} / 2} \leq \frac{K_{5} \mu}{\gamma \omega} \exp \left(-\chi^{n}\right) \tag{41}
\end{equation*}
$$

Proof. Let $\varphi: \mathbb{R}^{2} \rightarrow \mathbb{R}^{+}$be a $\mathcal{C}^{\infty}$ function supported in the open ball $B(0,1)$ of center 0 and radius 1 and with $\int_{\mathbb{R}^{2}} \varphi=1$. Let $\varphi_{n}: \mathbb{R}^{2} \rightarrow \mathbb{R}^{+}$be the mollifier

$$
\varphi_{n}(x):=\frac{N_{n}^{6}}{\nu_{0}^{2} \gamma^{8}} \varphi\left(\frac{N_{n}^{3}}{\nu_{0} \gamma^{4}} x\right)
$$

$\operatorname{Supp}\left(\varphi_{n}\right) \subset B\left(0, \nu_{0} \gamma^{4} / N_{n}^{3}\right)$ and $\int_{\mathbb{R}^{2}} \varphi_{n}=1$. We define $\psi_{n}: \mathbb{R}^{2} \rightarrow \mathbb{R}$ as

$$
\psi_{n}(x):=\left(\varphi_{n} * \chi_{A_{n}^{*}}\right)(x)=\int_{\mathbb{R}^{2}} \varphi_{n}(y-x) \chi_{A_{n}^{*}}(y) d y
$$

where $\chi_{A_{n}^{*}}$ is the characteristic function of the set $A_{n}^{*} . \psi_{n}$ is $\mathcal{C}^{\infty}$,

$$
\begin{equation*}
\left|D \psi_{n}(x)\right| \leq \int_{\mathbb{R}^{2}}\left|D \varphi_{n}(x-y)\right| \chi_{A_{n}^{*}}(y) d y \leq \frac{N_{n}^{3}}{\nu_{0} \gamma^{4}} C \tag{42}
\end{equation*}
$$

where $C:=\int_{\mathbb{R}^{2}}|D \varphi| d y$,

$$
0 \leq \psi_{n}(x) \leq 1, \quad \operatorname{supp}\left(\psi_{n}\right) \subset A_{n}, \quad \psi_{n}(x)=1 \quad \forall x \in \tilde{A}_{n}
$$

We define, for $(\mu, \omega) \in A_{0}$, the $C^{\infty}$ functions

$$
\tilde{h}_{n}(\mu, \omega):= \begin{cases}\psi_{n}(\mu, \omega) h_{n}(\mu, \omega) & \text { if } \quad(\mu, \omega) \in A_{n} \\ 0 & \text { if } \quad(\mu, \omega) \notin A_{n}\end{cases}
$$

and

$$
\tilde{w}_{n}(\mu, \omega):=\sum_{i=0}^{n} \tilde{h}_{i}, \quad \tilde{w}(\mu, \omega):=\sum_{i \geq 0} \tilde{h}_{i}
$$

which is a series if $(\mu, \omega) \in A_{\infty}:=\bigcap_{n \geq 0} A_{n}$.

The estimate for $\|\tilde{w}\|_{\sigma_{0} / 2}$ follows by $\left\|\tilde{h}_{i}\right\|_{\sigma_{i}} \leq\left\|h_{i}\right\|_{\sigma_{i}}$ (because $0 \leq \psi_{i} \leq 1$ ) and (28). The estimates for the derivatives in (40) follow differentiating the product $\tilde{h}_{i}=\psi_{i} h_{i}$ and using (42), (28) and Lemma 9. Similarly it follows that $\tilde{w}$ is in $C^{\infty}$, see [7] for details.

For $(\mu, \omega) \in \tilde{A}_{n}, \psi_{n}(\mu, \omega)=1$, implying $\tilde{w}_{n}=w_{n}$. As a consequence, for $(\mu, \omega) \in \tilde{A}_{\infty}:=\bigcap_{n>0} \tilde{A}_{n}$, by Corollary 1, $\tilde{w}=w_{\infty}$ solves (19).

Finally, using (28),

$$
\left\|\tilde{w}-\tilde{w}_{n}\right\|_{\sigma_{0} / 2} \leq \sum_{i \geq n+1}\left\|\tilde{h}_{i}\right\|_{\sigma_{i}} \leq \sum_{i \geq n+1} \frac{K \mu}{\gamma \omega} \exp \left(-\chi^{i}\right) \leq \frac{K^{\prime} \mu}{\gamma \omega} \exp \left(-\chi^{n}\right)
$$

Lemma 11. There exist $K_{5}^{\prime}$ such that if $\mu / \gamma^{2} \omega<K_{5}^{\prime}$ and $\nu_{0}<K_{5}^{\prime}$ then

$$
B_{\gamma} \subseteq \tilde{A}_{n} \quad \forall n \geq 0
$$

where $B_{\gamma}$ is defined in (11).
Proof. By induction. Let $(\mu, \omega) \in B_{\gamma}$. Then $(\mu, \omega) \in \tilde{A}_{0}$ if $A_{0}$ contains the closed ball of center $(\mu, \omega)$ and radius $2 \nu_{0} \gamma^{4} / N_{0}^{3}$. Let ( $\omega^{\prime}, \mu^{\prime}$ ) belong to such a ball. Then, $\forall l=1, \ldots, N_{0}$,

$$
\left|\omega^{\prime} l-\omega_{j}\right| \geq\left|\omega l-\omega_{j}\right|-\left|\omega-\omega^{\prime}\right| l>\frac{2 \gamma}{l^{\tau}}-\frac{2 \nu_{0} \gamma^{4}}{N_{0}^{3}} l \geq \frac{\gamma}{l^{\tau}}
$$

if $\nu_{0} \leq 1 / 2$.
Suppose now $B_{\gamma} \subseteq \tilde{A}_{n}$ and let $(\mu, \omega) \in B_{\gamma}$. To prove that $(\mu, \omega) \in \tilde{A}_{n+1}$, we have to show that the closed ball of center $(\mu, \omega)$ and radius $2 \nu_{0} \gamma^{4} / N_{n+1}^{3}$ is contained in $A_{n+1}$. Let $\left(\mu^{\prime}, \omega^{\prime}\right)$ belong to such a ball. The non-resonance condition on $\left|\omega^{\prime} l-j / c\right|$ is verified, as above, for $\nu_{0} \leq 1 / 2$. For the other condition, we denote in short $\omega_{j}^{n}\left(\mu^{\prime}, \omega^{\prime}\right):=\omega_{j}\left(\mu^{\prime}, w_{n}\left(\mu^{\prime}, \omega^{\prime}\right)\right)$ (see (16) for the definition of $\left.\omega_{j}(\mu, w)\right)$. It results, $\forall l=1, \ldots, N_{n+1}$,

$$
\begin{align*}
\left|\omega^{\prime} l-\omega_{j}^{n}\left(\mu^{\prime}, \omega^{\prime}\right)\right| & \geq\left|\omega l-\tilde{\omega}_{j}(\mu, \omega)\right|-\left|\omega-\omega^{\prime}\right| l-\left|\omega_{j}^{n}\left(\mu^{\prime}, \omega^{\prime}\right)-\tilde{\omega}_{j}(\mu, \omega)\right| \\
& >\frac{2 \gamma}{l^{\tau}}-\frac{2 \nu_{0} \gamma^{4} l}{N_{n+1}^{3}}-\left|\omega_{j}^{n}\left(\mu^{\prime}, \omega^{\prime}\right)-\tilde{\omega}_{j}(\mu, \omega)\right| \\
& >\frac{3 \gamma}{2 l^{\tau}}-\left|\omega_{j}^{n}\left(\mu^{\prime}, \omega^{\prime}\right)-\tilde{\omega}_{j}(\mu, \omega)\right| \tag{43}
\end{align*}
$$

if $\nu_{0} \leq 1 / 4$. We now estimate the last term

$$
\left|\omega_{j}^{n}\left(\mu^{\prime}, \omega^{\prime}\right)-\tilde{\omega}_{j}(\mu, \omega)\right|=\frac{\left|\lambda_{j}^{n}\left(\mu^{\prime}, \omega^{\prime}\right)-\tilde{\lambda}_{j}(\mu, \omega)\right|}{\left|\tilde{\omega}_{j}(\mu, \omega)\right|+\left|\omega_{j}^{n}\left(\mu^{\prime}, \omega^{\prime}\right)\right|} \leq \frac{\left|\lambda_{j}^{n}\left(\mu^{\prime}, \omega^{\prime}\right)-\tilde{\lambda}_{j}(\mu, \omega)\right|}{\sqrt{\delta_{0}}}
$$

by (18), both for $j<j_{0}$ and for $j \geq j_{0}$. By the comparison principle (17)

$$
\delta_{0}^{-1 / 2}\left|\lambda_{j}^{n}\left(\mu^{\prime}, \omega^{\prime}\right)-\tilde{\lambda}_{j}(\mu, \omega)\right| \leq K\left|\mu-\mu^{\prime}\right|+K\left\|w_{n}\left(\mu^{\prime}, \omega^{\prime}\right)-\tilde{w}(\mu, \omega)\right\|_{\sigma_{0} / 2} .
$$

By Lemma $9,\left\|\partial_{\omega} w_{n}\right\|_{\sigma_{0} / 2},\left\|\partial_{\mu} w_{n}\right\|_{\sigma_{0} / 2} \leq K_{0} / \omega \gamma^{2}$, and being $\omega, \omega^{\prime}>\gamma$,

$$
K\left\|w_{n}\left(\mu^{\prime}, \omega^{\prime}\right)-w_{n}(\mu, \omega)\right\|_{\sigma_{0} / 2} \leq \frac{K^{\prime}}{\gamma^{3}} \frac{\nu_{0} \gamma^{4}}{N_{n+1}^{3}}<\frac{\gamma}{8 l^{\tau}}, \forall l=1, \ldots, N_{n+1}
$$

if $\nu_{0}$ is small enough $(1<\tau<2)$. On the other hand, since $(\mu, \omega) \in \tilde{A}_{n}$ we have $w_{n}(\mu, \omega)=\tilde{w}_{n}(\mu, \omega)($ Lemma 10) and, by (41),

$$
K\left\|w_{n}(\mu, \omega)-\tilde{w}(\mu, \omega)\right\|_{\sigma_{0} / 2} \leq \frac{K^{\prime \prime} \mu}{\gamma \omega} \exp \left(-\chi^{n}\right)<\frac{\gamma}{8 l^{\tau}}, \forall l=1, \ldots, N_{n+1}
$$

for $\mu / \gamma^{2} \omega$ sufficiently small. By (43), collecting the previous estimates,

$$
\left|\omega^{\prime} l-\omega_{j}^{n}\left(\mu^{\prime}, \omega^{\prime}\right)\right|>\frac{\gamma}{l^{\tau}}, \quad \forall l=1, \ldots, N_{n+1}
$$

and ( $\mu^{\prime}, \omega^{\prime}$ ) belongs to $A_{n+1}$.

### 3.3 Measure of the Cantor set $B_{\gamma}$

In the following $R:=\left(\mu^{\prime}, \mu^{\prime \prime}\right) \times\left(\omega^{\prime}, \omega^{\prime \prime}\right)$ denotes a rectangle contained in the region $\left\{(\mu, \omega) \in\left[\mu_{1}, \mu_{2}\right] \times(2 \gamma,+\infty): \mu<K_{6}^{\prime} \gamma^{5} \omega\right\}$. Furthermore we consider $\omega^{\prime \prime}-\omega^{\prime}$ as a fixed quantity ("of order 1 ").

Lemma 12. There exist $K_{6}, K_{6}^{\prime}$ such that, $\forall \mu \in\left[\mu_{1}, \mu_{2}\right]$, the section

$$
S_{\gamma}(\mu):=\left\{\omega:(\mu, \omega) \in B_{\gamma}\right\}
$$

with $\mu / \omega \gamma^{5}<K_{6}^{\prime}$ in the definition (11) of $B_{\gamma}$, satisfies the measure estimate

$$
\begin{equation*}
\left|S_{\gamma}(\mu) \cap\left(\omega^{\prime}, \omega^{\prime \prime}\right)\right| \geq\left(1-K_{6} \gamma\right)\left(\omega^{\prime \prime}-\omega^{\prime}\right) \tag{44}
\end{equation*}
$$

for $\gamma$ small. As a consequence, for every $R:=\left(\mu^{\prime}, \mu^{\prime \prime}\right) \times\left(\omega^{\prime}, \omega^{\prime \prime}\right)$

$$
\begin{equation*}
\left|B_{\gamma} \cap R\right| \geq|R|\left(1-K_{6} \gamma\right) \tag{45}
\end{equation*}
$$

Proof. We consider just the inequalities $\left|\omega l-\tilde{\omega}_{j}(\mu, \omega)\right|>2 \gamma / l^{\tau}$ in the definition of $B_{\gamma}$. The analogous inequalities are simpler because $j / c$ and $\omega_{j}$ do not depend on $(\mu, \omega)$.

The complementary set we have to estimate is

$$
\mathcal{C}:=\bigcup_{l, j \geq 1} \mathcal{R}_{l j}
$$

where $\mathcal{R}_{l j}:=\left\{\omega \in\left(\omega^{\prime}, \omega^{\prime \prime}\right):\left|l \omega-\tilde{\omega}_{j}(\mu, \omega)\right| \leq 2 \gamma / l^{\tau}\right\}$.
We claim that

$$
\begin{equation*}
\left|\partial_{\omega} \tilde{\omega}_{j}(\mu, \omega)\right| \leq \frac{K \mu}{\gamma^{5} \omega} . \tag{46}
\end{equation*}
$$

Indeed, by the same arguments as in the proof of Lemma 11 and the comparison principle (17) we have

$$
\left|\tilde{\omega}_{j}(\mu, \omega)-\tilde{\omega}_{j}\left(\mu, \omega^{\prime}\right)\right| \leq K\left\|\tilde{w}(\mu, \omega)-\tilde{w}\left(\mu, \omega^{\prime}\right)\right\|_{\sigma_{0} / 2} \leq \frac{K \mu}{\gamma^{5} \omega}\left|\omega-\omega^{\prime}\right|
$$

using (40). As a consequence of (46)

$$
\partial_{\omega}\left(l \omega-\tilde{\omega}_{j}(\mu, \omega)\right) \geq l-\frac{K \mu}{\gamma^{5} \omega} \geq \frac{l}{2} \quad \forall l \geq 1
$$

for $\mu / \gamma^{5} \omega$ small enough; we deduce $\left|\mathcal{R}_{l j}\right| \leq 4 \gamma / l^{\tau+1}$.
Furthermore the set $\mathcal{R}_{l j}$ is non-empty only if

$$
\omega^{\prime} l-\frac{2 \gamma}{l^{\tau}}<\tilde{\omega}_{j}(\mu, \omega)<\omega^{\prime \prime} l+\frac{2 \gamma}{l^{\tau}} .
$$

So, for every fixed $l$, the number of indices $j$ such that $\mathcal{R}_{l, j} \neq \emptyset$ is

$$
\sharp\{j\} \leq \frac{1}{\delta}\left(l\left(\omega^{\prime \prime}-\omega^{\prime}\right)+\frac{4 \gamma}{l^{\tau}}\right)+1 \leq K l\left(\omega^{\prime \prime}-\omega^{\prime}\right)
$$

where

$$
\delta:=\inf \left\{\left|\tilde{\omega}_{j+1}(\mu, \omega)-\tilde{\omega}_{j}(\mu, \omega)\right|: j \geq 1,(\mu, \omega) \in B_{\gamma}\right\}
$$

For $\|\tilde{w}\|_{\sigma_{0} / 2} \leq K^{\prime} \mu / \gamma \omega<R$ we have $\delta \geq \delta_{1}$ where

$$
\begin{equation*}
\delta_{1}:=\inf \left\{\left|\omega_{j+1}(\mu, w)-\omega_{j}(\mu, w)\right|: j \geq 1, \mu \in\left[\mu_{1}, \mu_{2}\right],\|w\|_{\sigma_{0} / 2} \leq R\right\}>0 \tag{47}
\end{equation*}
$$

as proved in the Appendix.
In conclusion, the measure of the complementary set is

$$
|\mathcal{C}| \leq \sum_{l=1}^{+\infty} \frac{4 \gamma}{l^{\tau+1}} K l\left(\omega^{\prime \prime}-\omega^{\prime}\right) \leq K^{\prime}\left(\omega^{\prime \prime}-\omega^{\prime}\right) \gamma
$$

and (44) is proved. Integrating on ( $\mu^{\prime}, \mu^{\prime \prime}$ ) we obtain (45).
By Fubini Theorem also the section $S_{\gamma}(\omega)$ is large, for $\omega$ in a large set.

Lemma 13. Let

$$
S_{\gamma}(\omega):=\left\{\mu:(\mu, \omega) \in B_{\gamma}\right\} .
$$

For every $R:=\left(\mu^{\prime}, \mu^{\prime \prime}\right) \times\left(\omega^{\prime}, \omega^{\prime \prime}\right), \gamma^{\prime} \in(0,1)$ it results

$$
\begin{equation*}
\left|\left\{\omega \in\left(\omega^{\prime}, \omega^{\prime \prime}\right): \frac{\left|S_{\gamma}(\omega) \cap\left(\mu^{\prime}, \mu^{\prime \prime}\right)\right|}{\mu^{\prime \prime}-\mu^{\prime}} \geq 1-\gamma^{\prime}\right\}\right| \geq\left(\omega^{\prime \prime}-\omega^{\prime}\right)\left(1-K_{6} \frac{\gamma}{\gamma^{\prime}}\right) . \tag{48}
\end{equation*}
$$

Proof. Let consider

$$
\begin{aligned}
& \Omega^{+}:=\left\{\omega \in\left(\omega^{\prime}, \omega^{\prime \prime}\right):\left|S_{\gamma}(\omega) \cap\left(\mu^{\prime}, \mu^{\prime \prime}\right)\right| \geq\left(\mu^{\prime \prime}-\mu^{\prime}\right)\left(1-\gamma^{\prime}\right)\right\} \\
& \Omega^{-}:=\left\{\omega \in\left(\omega^{\prime}, \omega^{\prime \prime}\right):\left|S_{\gamma}(\omega) \cap\left(\mu^{\prime}, \mu^{\prime \prime}\right)\right|<\left(\mu^{\prime \prime}-\mu^{\prime}\right)\left(1-\gamma^{\prime}\right)\right\} .
\end{aligned}
$$

Using the Fubini theorem

$$
\begin{align*}
\left|B_{\gamma} \cap R\right| & =\int_{\omega^{\prime}}^{\omega^{\prime \prime}}\left|S_{\gamma}(\omega) \cap\left(\mu^{\prime}, \mu^{\prime \prime}\right)\right| d \omega  \tag{49}\\
& =\int_{\Omega^{+}}\left|S_{\gamma}(\omega) \cap\left(\mu^{\prime}, \mu^{\prime \prime}\right)\right| d \omega+\int_{\Omega^{-}}\left|S_{\gamma}(\omega) \cap\left(\mu^{\prime}, \mu^{\prime \prime}\right)\right| d \omega \\
& \leq\left(\mu^{\prime \prime}-\mu^{\prime}\right)\left|\Omega^{+}\right|+\left(\mu^{\prime \prime}-\mu^{\prime}\right)\left(1-\gamma^{\prime}\right)\left|\Omega^{-}\right| .
\end{align*}
$$

Minorating the left hand side in (49) by (45) yields

$$
\begin{equation*}
\left(\omega^{\prime \prime}-\omega^{\prime}\right)\left(1-K_{6} \gamma\right) \leq\left|\Omega^{+}\right|+\left(1-\gamma^{\prime}\right)\left|\Omega^{-}\right|=\left(\omega^{\prime \prime}-\omega^{\prime}\right)-\gamma^{\prime}\left|\Omega^{-}\right| \tag{50}
\end{equation*}
$$

and therefore $\left|\Omega^{-}\right| \leq\left(\omega^{\prime \prime}-\omega^{\prime}\right) K_{6} \gamma / \gamma^{\prime}$. We deduce by the first inequality in (50) that $\left|\Omega^{+}\right| \geq\left(\omega^{\prime \prime}-\omega^{\prime}\right)\left(1-K_{6} \gamma / \gamma^{\prime}\right)$, namely (48).

By (44)-(48) we deduce the measure estimates for the "sections" (in $\omega$ and $\mu$ ) of $\cup_{\gamma \in(0,1)} B_{\gamma}$ stated after Theorem 1.

## 4 Inversion of the linearized problem

Here we prove Lemma 5. Decomposing in Fourier series $f^{\prime}(u)=\sum_{k \in \mathbb{Z}} a_{k}(x)$ $e^{i k t}$ we have, $\forall h=\sum_{1 \leq|l| \leq N_{n}} h_{l}(x) e^{i l t} \in W^{(n)}$,

$$
\begin{align*}
-L_{\omega} h+\mu P_{n}\left[f^{\prime}(u) h\right]= & \sum_{1 \leq|l| \leq N_{n}}\left[\omega^{2} l^{2} \rho h_{l}+\partial_{x}\left(p \partial_{x} h_{l}\right)\right] e^{i l t}+ \\
& +\mu P_{n}\left[\left(\sum_{k \in \mathbb{Z}} a_{k} e^{i k t}\right)\left(\sum_{1 \leq|l| \leq N_{n}} h_{l} e^{i l t}\right)\right] \\
= & \sum_{1 \leq|l| \leq N_{n}}\left[\omega^{2} l^{2} \rho h_{l}+\partial_{x}\left(p \partial_{x} h_{l}\right)+\mu a_{0} h_{l}\right] e^{i l t}  \tag{51}\\
& +\mu \sum_{\left||l|,|k+l| \in\left\{1, \ldots, N_{n}\right\}, k \neq 0\right.} a_{k} h_{l} e^{i(k+l) t} \tag{52}
\end{align*}
$$

distinguishing the diagonal operator (51) by the off-diagonal term (52). Hence $\mathcal{L}_{n}(w)$ defined in (21) can be decomposed as

$$
\begin{equation*}
\mathcal{L}_{n}(w) h=\rho\left(D h+M_{1} h+M_{2} h\right) \tag{53}
\end{equation*}
$$

where

$$
\begin{align*}
D h & :=\frac{1}{\rho} \sum_{|l|=1}^{N_{n}}\left[\omega^{2} l^{2} \rho h_{l}+\left(p h_{l}^{\prime}\right)^{\prime}+\mu a_{0} h_{l}\right] e^{i l t} \\
M_{1} h & :=\frac{\mu}{\rho} \sum_{|l|,|k| \in\left\{1, \ldots, N_{n}\right\}, l \neq k} a_{k-l} h_{l} e^{i k t}  \tag{54}\\
M_{2} h & :=\frac{\mu}{\rho} P_{n}\left[f^{\prime}(u) d_{w} v(\mu, w)[h]\right] .
\end{align*}
$$

To study the eigenvalues of $D$, we use Sturm-Liouville type techniques.
Lemma 14. (Sturm-Liouville) The eigenvalues $\lambda_{j}(\mu, w)$ of the SturmLiouville problem (16) form a strictly increasing sequence which tends to $+\infty$. Every $\lambda_{j}(\mu, w)$ is simple and the following asymptotic formula holds

$$
\begin{equation*}
\lambda_{j}(\mu, w)=\frac{j^{2}}{c^{2}}+b+M(\mu, w)+r_{j}(\mu, w), \quad\left|r_{j}(\mu, w)\right| \leq \frac{K}{j} \tag{55}
\end{equation*}
$$

$\forall j \geq 1,(\mu, w) \in\left[\mu_{1}, \mu_{2}\right] \times B_{R}$, where

$$
\begin{aligned}
& c:=\frac{1}{\pi} \int_{0}^{\pi}\left(\frac{\rho}{p}\right)^{1 / 2} d x, \quad b:=\frac{1}{4 \pi c} \int_{0}^{\pi}\left[\frac{(\rho p)^{\prime}}{\rho \sqrt[4]{\rho p}}\right]^{\prime} \frac{1}{\sqrt[4]{\rho p}} d x \\
& M(\mu, w):=-\frac{\mu}{c \pi} \int_{0}^{\pi} \frac{\Pi_{V} f^{\prime}(v(\mu, w)+w)}{\sqrt{\rho p}} d x
\end{aligned}
$$

The eigenfunctions $\varphi_{j}(\mu, w)$ of (16) form an orthonormal basis of $L^{2}(0, \pi)$ with respect to the scalar product $(y, z)_{L_{\rho}^{2}}:=c^{-1} \int_{0}^{\pi} y z \rho d x$. For $K$ big enough

$$
(y, z)_{\mu, w}:=\frac{1}{c} \int_{0}^{\pi} p y^{\prime} z^{\prime}+\left[K \rho-\mu \Pi_{V} f^{\prime}(v(\mu, w)+w)\right] y z d x
$$

defines an equivalent scalar product on $H_{0}^{1}(0, \pi)$ and

$$
\begin{equation*}
K^{\prime}\|y\|_{H^{1}} \leq\|y\|_{\mu, w} \leq K^{\prime \prime}\|y\|_{H^{1}} \quad \forall y \in H_{0}^{1} \tag{56}
\end{equation*}
$$

$\varphi_{j}(\mu, w)$ is also an orthogonal basis of $H_{0}^{1}(0, \pi)$ with respect to the scalar product $(,)_{\mu, w}$ and, for $y=\sum_{j \geq 1} \hat{y}_{j} \varphi_{j}(\mu, w)$,

$$
\begin{equation*}
\|y\|_{L_{\rho}^{2}}^{2}=\sum_{j \geq 1} \hat{y}_{j}^{2}, \quad\|y\|_{\mu, w}^{2}=\sum_{j \geq 1} \hat{y}_{j}^{2}\left(\lambda_{j}(\mu, w)+K\right) \tag{57}
\end{equation*}
$$

Proof. In the Appendix.
We develop $D h=\sum_{1 \leq|l| \leq N_{n}} D_{l} h_{l} e^{i l t}$ where

$$
D_{l} z:=\frac{1}{\rho}\left[\omega^{2} l^{2} \rho z+\left(p z^{\prime}\right)^{\prime}+\mu a_{0} z\right], \quad \forall z \in H_{0}^{1}(0, \pi)
$$

and $a_{0}=\Pi_{V} f(v(\mu, w)+w)$.
By Lemma 14 each $D_{l}$ is diagonal w.r.t the basis $\varphi_{j}(\mu, w)$ :
$z=\sum_{j=1}^{+\infty} \hat{z}_{j} \varphi_{j}(\mu, w) \in H_{0}^{1}(0, \pi) \Rightarrow D_{l} z=\sum_{j=1}^{+\infty}\left(\omega^{2} l^{2}-\lambda_{j}(\mu, w)\right) \hat{z}_{j} \varphi_{j}(\mu, w)$.
Lemma 15. Suppose all the eigenvalues $\omega^{2} l^{2}-\lambda_{j}(\mu, w)$ are not zero. Then

$$
\left|D_{l}\right|^{-1 / 2} z:=\sum_{j=1}^{+\infty} \frac{\hat{z}_{j} \varphi_{j}(\mu, w)}{\sqrt{\left|\omega^{2} l^{2}-\lambda_{j}(\mu, w)\right|}}
$$

satisfies

$$
\begin{equation*}
\left\|\left|D_{l}\right|^{-1 / 2} z\right\|_{H^{1}} \leq \frac{K}{\sqrt{\alpha_{l}}}\|z\|_{H^{1}}, \quad \forall z \in H_{0}^{1}(0, \pi) \tag{58}
\end{equation*}
$$

where $\alpha_{l}:=\min _{j \geq 1}\left|\omega^{2} l^{2}-\lambda_{j}(\mu, w)\right|>0$.
Proof. By (57) $\left\|\left|D_{l}\right|^{-1 / 2} z\right\|_{\mu, w}^{2} \leq\left(1 / \alpha_{l}\right)\|z\|_{\mu, w}^{2}$. Hence (58) follows by the equivalence of the norms (56).

Lemma 16. (Inversion of $D$ ) Assume the non-resonance condition (23). Then $|D|^{-1 / 2}: W^{(n)} \rightarrow W^{(n)}$ defined by

$$
|D|^{-1 / 2} h:=\sum_{1 \leq|l| \leq N_{n}}\left|D_{l}\right|^{-1 / 2} h_{l} e^{i l t}
$$

satisfies

$$
\left\||D|^{-1 / 2} h\right\|_{\sigma, s} \leq \frac{K}{\sqrt{\gamma \omega}}\|h\|_{\sigma, s+\frac{\tau-1}{2}} \leq \frac{K N_{n}^{\frac{\tau-1}{2}}}{\sqrt{\gamma \omega}}\|h\|_{\sigma, s}, \quad \forall h \in W^{(n)}
$$

Proof. By (58) and $\alpha_{-l}=\alpha_{l} \geq \gamma \omega /|l|^{\tau-1}$

$$
\begin{aligned}
\left\||D|^{-1 / 2} h\right\|_{\sigma, s}^{2} & =\sum_{1 \leq|l| \leq N_{n}}\left\|\left|D_{l}\right|^{-1 / 2} h_{l}\right\|_{H^{1}}^{2}\left(1+l^{2 s}\right) e^{2 \sigma|l|} \\
& \leq \sum_{1 \leq|l| \leq N_{n}} \frac{K^{2}|l|^{\tau-1}}{\gamma \omega}\left\|h_{l}\right\|_{H^{1}}^{2}\left(1+l^{2 s}\right) e^{2 \sigma|l|} \\
& \leq \frac{K^{\prime}}{\gamma \omega}\|h\|_{\sigma, s+\frac{\tau-1}{2}}^{2}
\end{aligned}
$$

because $|l|^{\tau-1}\left(1+l^{2 s}\right)<2\left(1+|l|^{2 s+\tau-1}\right), \forall|l| \geq 1$.
To prove the invertibility of $\mathcal{L}_{n}(w)$ we write (53) as

$$
\begin{equation*}
\mathcal{L}_{n}(w)=\rho|D|^{1 / 2}\left(U+T_{1}+T_{2}\right)|D|^{1 / 2} \tag{59}
\end{equation*}
$$

where

$$
\left\{\begin{array}{l}
U:=|D|^{-1 / 2} D|D|^{-1 / 2}  \tag{60}\\
T_{i}:=|D|^{-1 / 2} M_{i}|D|^{-1 / 2}, \quad i=1,2
\end{array}\right.
$$

With respect to the basis $\varphi_{j}(\mu, w) e^{i l t}$ the operator $U$ is diagonal and its $(l, j)$-th eigenvalue is $\operatorname{sign}\left(\omega^{2} l^{2}-\lambda_{j}(\mu, w)\right) \in\{ \pm 1\}$, implying ${ }^{5}\|U\|_{\sigma}=1$.

The smallness of $T_{1}$ requires an analysis of the small divisors. Formula (55) implies, by Taylor expansion, the asymptotic dispersion relation

$$
\begin{equation*}
\left|\omega_{j}(\mu, w)-\frac{j}{c}\right| \leq \frac{K}{j} \tag{61}
\end{equation*}
$$

and there exists $K$ such that, for every $x \geq 0$,

$$
\begin{equation*}
\left|x^{2}-\lambda_{j^{*}}(\mu, w)\right|=\min _{j \geq 1}\left|x^{2}-\lambda_{j}(\mu, w)\right| \quad \Rightarrow \quad j^{*} \geq K x . \tag{62}
\end{equation*}
$$

Lemma 17. Assume the non-resonance conditions (22)-(23) and $\omega>\gamma$. Then $\forall|k|,|l| \in\left\{1, \ldots, N_{n}\right\}, k \neq l$

$$
\alpha_{l} \alpha_{k} \geq\left(\frac{K \gamma^{3} \omega}{|k-l|^{\frac{\tau(\tau-1)}{2-\tau}}}\right)^{2}
$$

where $\alpha_{l}:=\min _{j \geq 1}\left|\omega^{2} l^{2}-\lambda_{j}(\mu, w)\right|$.
Proof. Since $\alpha_{-l}=\alpha_{l}, \forall l$, we can suppose $l, k \geq 1$.
We distinguish two cases, if $k, l$ are close or far one from each other. Let $\beta:=(2-\tau) / \tau \in(0,1)$.
Case 1. Let $2|k-l|>(\max \{k, l\})^{\beta}$. By (23)

$$
\alpha_{k} \alpha_{l} \geq \frac{(\gamma \omega)^{2}}{(k l)^{\tau-1}} \geq \frac{(\gamma \omega)^{2}}{(\max \{k, l\})^{2(\tau-1)}} \geq \frac{C(\gamma \omega)^{2}}{|k-l|^{\frac{2(\tau-1)}{\beta}}} .
$$

Case 2. Let $0<2|k-l| \leq(\max \{k, l\})^{\beta}$. In this case $2 k \geq l \geq k / 2$. Indeed, if

[^4]$k>l$, then $2(k-l) \leq k^{\beta}$, so $2 l \geq 2 k-k^{\beta} \geq k$ because $\beta \in(0,1)$. Analogously if $l>k$.

Let $i$, resp. $j$, be an integer which realizes the minimum $\alpha_{k}$, resp. $\alpha_{l}$, and write in short $\lambda_{j}(\mu):=\lambda_{j}(\mu, w), \omega_{j}(\mu):=\omega_{j}(\mu, w)$.

If both $\lambda_{i}(\mu), \lambda_{j}(\mu) \leq 0$, then $\alpha_{l} \geq \omega^{2} l^{2}, \alpha_{k} \geq \omega^{2} k^{2}, \alpha_{l} \alpha_{k} \geq \omega^{4}>\gamma^{2} \omega^{2}$.
If only $\lambda_{j}(\mu) \leq 0$, then $\alpha_{l} \alpha_{k} \geq \gamma \omega^{3} l^{2} / k^{\tau-1} \geq 2^{1-\tau} \gamma \omega^{3} \geq 2^{1-\tau} \gamma^{2} \omega^{2}$.
The really resonant cases happen if $\lambda_{i}(\mu), \lambda_{j}(\mu)>0$.
Suppose, for example, $\max \{k, l\}=k$. By (61), $\left|\omega_{j}(\mu)-(j / c)\right| \leq K / j$, and, by (62), $i \geq K \omega k, j \geq K \omega l$. Hence, using also (22),

$$
\begin{aligned}
\left|\left(\omega k-\omega_{i}(\mu)\right)-\left(\omega l-\omega_{j}(\mu)\right)\right| & =\left|\omega(k-l)-\left(\omega_{i}(\mu)-\omega_{j}(\mu)\right)\right| \\
& \geq\left|\omega(k-l)-\frac{i-j}{c}\right|-\frac{K}{\omega l}-\frac{K}{\omega k} \\
& \geq \frac{\gamma}{(k-l)^{\tau}}-\frac{3 K}{\omega k} \geq \frac{2^{\tau} \gamma}{k^{\beta \tau}}-\frac{3 K}{\omega k}
\end{aligned}
$$

because $2(k-l) \leq k^{\beta}, 2 l \geq k$. Since $\beta \tau<1$ and $k \leq 2 l$,

$$
\left|\left(\omega k-\omega_{i}(\mu)\right)-\left(\omega l-\omega_{j}(\mu)\right)\right| \geq \frac{1}{2}\left(\frac{\gamma}{k^{\beta \tau}}+\frac{\gamma}{l^{\beta \tau}}\right) \quad \forall k \geq\left(\frac{K}{\omega \gamma}\right)^{\frac{1}{1-\beta \tau}}:=k^{*} .
$$

The same conclusion if $\max \{k, l\}=l$. It follows that, for $\max \{k, l\} \geq$ $k^{*}$, there holds $\left|\omega k-\omega_{i}(\mu)\right| \geq \gamma / 2 k^{\beta \tau}$ or $\left|\omega l-\omega_{j}(\mu)\right| \geq \gamma / 2 l^{\beta \tau}$. Suppose $\left|\omega k-\omega_{i}(\mu)\right| \geq \gamma / 2 k^{\beta \tau}$. Then

$$
\alpha_{k}=\left|\omega^{2} k^{2}-\omega_{i}^{2}(\mu)\right| \geq\left|\omega k-\omega_{i}(\mu)\right| \omega k \geq \frac{\gamma \omega}{2} k^{1-\beta \tau} .
$$

Since $l \leq 2 k$, for $\alpha_{l}$ we can use (23),

$$
\alpha_{k} \alpha_{l} \geq \frac{\gamma \omega k^{1-\beta \tau}}{2} \frac{\gamma \omega}{l^{\tau-1}} \geq \frac{\gamma^{2} \omega^{2}}{2^{\tau}} k^{2-\tau-\beta \tau}=\frac{\gamma^{2} \omega^{2}}{2^{\tau}}
$$

because $2-\tau-\beta \tau=0$.
On the other hand, if $\max \{k, l\}<k^{*}=(K / \omega \gamma)^{1 /(\tau-1)}$, we can use (23) for both $k, l$ :

$$
\alpha_{k} \alpha_{l} \geq \frac{(\gamma \omega)^{2}}{(k l)^{\tau-1}}>\frac{(\gamma \omega)^{2}}{\left(k^{*}\right)^{2(\tau-1)}}=(\gamma \omega)^{2}\left(\frac{\omega \gamma}{K}\right)^{\frac{1}{\tau-1} 2(\tau-1)}>\frac{\gamma^{6} \omega^{2}}{K^{2}} .
$$

Since $\gamma<1$, taking the minimum for all these cases we conclude.

Lemma 18. (Estimate of $T_{1}$ ) Assume the non-resonance conditions (22)(23), $\omega>\gamma$, and $\Pi_{W} f^{\prime}(u)=\sum_{l \neq 0} a_{l}(x) e^{i l t} \in X_{\sigma, 1+\frac{\tau(\tau-1)}{2-\tau}}$. There exists $K$ such that

$$
\left\|T_{1} h\right\|_{\sigma} \leq \frac{K \mu}{\gamma^{3} \omega}\left\|\Pi_{W} f^{\prime}(u)\right\|_{\sigma, 1+\frac{\tau(\tau-1)}{2-\tau}}\|h\|_{\sigma}, \quad \forall h \in W^{(n)}
$$

Proof. $\forall h \in W^{(n)}, T_{1} h=\sum_{1 \leq|k| \leq N_{n}}\left(T_{1} h\right)_{k} e^{i k t}$ where

$$
\begin{aligned}
\left(T_{1} h\right)_{k} & =\left|D_{k}\right|^{-1 / 2}\left(M_{1}|D|^{-1 / 2} h\right)_{k} \\
& =\left|D_{k}\right|^{-1 / 2}\left[\sum_{1 \leq|l| \leq N_{n}, l \neq k} \mu \frac{a_{k-l}}{\rho}\left|D_{l}\right|^{-1 / 2} h_{l}\right] .
\end{aligned}
$$

Setting $A_{m}:=\left\|a_{m} / \rho\right\|_{H^{1}}$, using (58) and Lemma 17,

$$
\begin{equation*}
\left\|\left(T_{1} h\right)_{k}\right\|_{H^{1}} \leq K \mu \sum_{1 \leq|l| \leq N_{n}, l \neq k} \frac{A_{k-l}}{\sqrt{\alpha_{k}} \sqrt{\alpha_{l}}}\left\|h_{l}\right\|_{H^{1}} \leq \frac{K \mu}{\gamma^{3} \omega} S_{k} \tag{63}
\end{equation*}
$$

where

$$
S_{k}:=\sum_{|l| \leq N_{n}, l \neq k} A_{k-l}|k-l|^{\frac{\tau(\tau-1)}{2-\tau}}\left\|h_{l}\right\|_{H^{1}}
$$

By (63) we get, defining $S(t):=\sum_{|k|=1}^{N_{n}} S_{k} e^{i k t}$,

$$
\begin{aligned}
\left\|T_{1} h\right\|_{\sigma}^{2} & =\sum_{|k|=1}^{N_{n}}\left\|\left(T_{1} h\right)_{k}\right\|_{H^{1}}^{2}\left(1+k^{2}\right) e^{2 \sigma|k|} \\
& \leq\left(\frac{K \mu}{\gamma^{3} \omega}\right)^{2} \sum_{|k|=1}^{N_{n}} S_{k}^{2}\left(1+k^{2}\right) e^{2 \sigma|k|}=\left(\frac{K \mu}{\gamma^{3} \omega}\right)^{2}\|S\|_{\sigma}^{2}
\end{aligned}
$$

Since $S=P_{n}(\varphi \psi)$ with $\varphi(t):=\sum_{l \in \mathbb{Z}} A_{l}|l|^{\frac{\tau(\tau-1)}{2-\tau}} e^{i l t}$ and $\psi(t):=\sum_{|l|=1}^{N_{n}}\left\|h_{l}\right\|_{H^{1}} e^{i l t}$

$$
\left\|T_{1} h\right\|_{\sigma} \leq \frac{K \mu}{\gamma^{3} \omega}\|\varphi\|_{\sigma}\|\psi\|_{\sigma} \leq \frac{K \mu}{\gamma^{3} \omega}\left\|\Pi_{W} f^{\prime}(u)\right\|_{\sigma, 1+\frac{\tau(\tau-1)}{2-\tau}}\|h\|_{\sigma}
$$

because $\|\varphi\|_{\sigma} \leq 2\left\|\Pi_{W} f^{\prime}(u)\right\|_{\sigma, 1+\frac{\tau(\tau-1)}{2-\tau}}$ and $\|\psi\|_{\sigma}=\|h\|_{\sigma}$.
Lemma 19. (Estimate of $\left.T_{2}\right)$ Suppose that $\Pi_{W} f^{\prime}(u) \in X_{\sigma, 1+\frac{\tau-1}{2}}$. Then

$$
\left\|T_{2} h\right\|_{\sigma} \leq \frac{K \mu}{\gamma \omega}\left\|\Pi_{W} f^{\prime}(u)\right\|_{\sigma, 1+\frac{\tau-1}{2}}\|h\|_{\sigma}, \quad \forall h \in W^{(n)}
$$

for some $K$.

Proof. By the definitions (60),(54) and Lemma 16,

$$
\begin{aligned}
\left\|T_{2} h\right\|_{\sigma} & \leq \frac{K}{\sqrt{\gamma \omega}}\left\|M_{2}|D|^{-1 / 2} h\right\|_{\sigma, 1+\frac{\tau-1}{2}} \\
& \leq \frac{K^{\prime} \mu}{\sqrt{\gamma \omega}}\left\|\Pi_{W} f^{\prime}(u)\right\|_{\sigma, 1+\frac{\tau-1}{2}}\left\|d_{w} v(\mu, w)\left[|D|^{-1 / 2} h\right]\right\|_{\sigma, 1+\frac{\tau-1}{2}} \\
& =\frac{K^{\prime} \mu}{\sqrt{\gamma \omega}}\left\|\Pi_{W} f^{\prime}(u)\right\|_{\sigma, 1+\frac{\tau-1}{2}}\left\|d_{w} v(\mu, w)\left[|D|^{-1 / 2} h\right]\right\|_{H^{1}}
\end{aligned}
$$

because $d_{w} v(\mu, w)\left[|D|^{-1 / 2} h\right] \in V$. By Lemmas 4 and 16

$$
\left\|d_{w} v(\mu, w)\left[|D|^{-1 / 2} h\right]\right\|_{H^{1}} \leq K\left\||D|^{-1 / 2} h\right\|_{\sigma, 1-\frac{\tau-1}{2}} \leq \frac{K}{\sqrt{\gamma \omega}}\|h\|_{\sigma, 1}
$$

implying the thesis.
Proof of Lemma 5. $\|U\|_{\sigma}=1$. If $\left\|T_{1}+T_{2}\right\|_{\sigma}<1 / 2$, then by Neumann series $U+T_{1}+T_{2}$ is invertible in $\left(W^{(n)},\| \|_{\sigma}\right)$ and $\left\|\left(U+T_{1}+T_{2}\right)^{-1}\right\|_{\sigma}<2$. By Lemmas 18, 19, this condition is verified if we choose $K_{1}$ in (24) small enough. Hence, inverting (59)

$$
\mathcal{L}_{n}(w)^{-1} h=|D|^{-1 / 2}\left(U+T_{1}+T_{2}\right)^{-1}|D|^{-1 / 2}\left(\frac{h}{\rho}\right)
$$

which, using Lemma 16, yields (25).

## 5 Appendix

Proof of Lemma 14. Let $a(x) \in L^{2}(0, \pi)$. Under the "Liouville change of variable"

$$
\begin{equation*}
x=\psi(\xi) \Leftrightarrow \xi=g(x), \quad g(x):=\frac{1}{c} \int_{0}^{x}\left(\frac{\rho(s)}{p(s)}\right)^{1 / 2} d s \tag{64}
\end{equation*}
$$

we have that $(\lambda, y(x))$ satisfies

$$
\left\{\begin{array}{l}
-\left(p(x) y^{\prime}(x)\right)^{\prime}+a(x) y(x)=\lambda \rho(x) y(x)  \tag{65}\\
y(0)=y(\pi)=0
\end{array}\right.
$$

if and only if $(\nu, z(\xi))$ satisfies

$$
\left\{\begin{array}{l}
-z^{\prime \prime}(\xi)+[q(\xi)+\alpha(\xi)] z(\xi)=\nu z(\xi)  \tag{66}\\
z(0)=z(\pi)=0
\end{array}\right.
$$

where

$$
\begin{aligned}
& \nu=c^{2} \lambda, \quad r(x)=\sqrt[4]{p(x) \rho(x)}, \quad z(\xi)=y(\psi(\xi)) r(\psi(\xi)), \\
& \alpha(\xi)=c^{2} \frac{a(\psi(\xi))}{\rho(\psi(\xi))}, \quad q(\xi)=c^{2} Q(\psi(\xi)), \quad Q=\frac{p}{\rho} \frac{r^{\prime \prime}}{r}+\frac{1}{2}\left(\frac{p}{\rho}\right)^{\prime} \frac{r^{\prime}}{r} .
\end{aligned}
$$

By [20], Theorem 4 in Chapter 2, p.35, the eigenvalues of (66) form an increasing sequence $\nu_{j}$ satisfying the asympototics

$$
\nu_{j}=j^{2}+\frac{1}{\pi} \int_{0}^{\pi}(q+\alpha) d \xi-\frac{1}{\pi} \int_{0}^{\pi} \cos (2 j \xi)(q(\xi)+\alpha(\xi)) d \xi+r_{j}, \quad\left|r_{j}\right| \leq \frac{C}{j}
$$

where $C:=C\left(\|q+\alpha\|_{L^{2}}\right)$ is a positive constant. Moreover every $\nu_{j}$ is simple ([20], Theorem 2, p.30).

Since $p, \rho$ are positive and belong to $H^{3}$, if $a \in H^{1}$ then $q, \alpha \in H^{1}$. Integrating by parts $\left|\int_{0}^{\pi} \cos (2 j \xi)(q+\alpha) d \xi\right| \leq\|q+\alpha\|_{H^{1}} / j$ and so

$$
\nu_{j}=j^{2}+\frac{1}{\pi} \int_{0}^{\pi}(q+\alpha) d \xi+r_{j}^{\prime}, \quad\left|r_{j}^{\prime}\right| \leq \frac{C^{\prime}}{j}
$$

for some $C^{\prime}:=C^{\prime}\left(\|q+\alpha\|_{H^{1}}\right)$. Dividing by $c^{2}$ and using the inverse Liouville change of variable we obtain the formula for the eigenvalues $\lambda_{j}(a)$ of (65)

$$
\begin{equation*}
\lambda_{j}(a)=\frac{j^{2}}{c^{2}}+\frac{1}{\pi c} \int_{0}^{\pi} \frac{Q \sqrt{\rho}}{\sqrt{p}} d x+\frac{1}{\pi c} \int_{0}^{\pi} \frac{a}{\sqrt{\rho p}} d x+r_{j}(a), \quad\left|r_{j}(a)\right| \leq \frac{C}{j} \tag{67}
\end{equation*}
$$

for some $C\left(\rho, p,\|a\|_{H^{1}}\right)>0$. Formula (55) follows for $a(x)=-\mu \Pi_{V} f^{\prime}(v(\mu, w)$ $+w)(x)$ and some algebra.

By [20], Theorem 7 p.43, the eigenfunctions of (66) form an orthonormal basis for $L^{2}$. Applying in the integrals the Liouville change of variable, the eigenfunctions $\varphi_{j}(a)$ of (65) form an orthonormal basis for $L^{2}$ w.r.t. the scalar product $(,)_{L_{\rho}^{2}}$.

Finally, since $\varphi_{j}:=\varphi_{j}(a)$ solves

$$
-\left(p \varphi_{j}^{\prime}\right)^{\prime}+(K \rho+a) \varphi_{j}=\left(\lambda_{j}(a)+K\right) \rho \varphi_{j},
$$

multiplying by $\varphi_{i}$ and integrating by parts gives

$$
\left(\varphi_{j}, \varphi_{i}\right)_{\mu, w}=\delta_{i, j}\left(\lambda_{j}(a)+K\right)
$$

and (57) follows (note that $\lambda_{j}(a)+K>0, \forall j$, for $K$ large enough).
Proof of (17). Let $a, b \in H^{1}(0, \pi)$ and consider $\alpha:=c^{2} a(\psi) / \rho(\psi)$, $\beta:=c^{2} b(\psi) / \rho(\psi)$ constructed as above via the Liouville change of variable (64). By [20], p.34, for every $j$

$$
\begin{equation*}
\left|\lambda_{j}(a)-\lambda_{j}(b)\right|=\frac{1}{c^{2}}\left|\nu_{j}(\alpha)-\nu_{j}(\beta)\right| \leq \frac{1}{c^{2}}\|\alpha-\beta\|_{\infty} \leq K\|a-b\|_{H^{1}} \tag{68}
\end{equation*}
$$

and (17) follows by the mean value theorem because $\mu \Pi_{V} f(v(\mu, w)+w)$ has bounded derivatives on bounded sets.

Proof of (47). By the asymptotic formula (61)

$$
\min _{j \geq 1}\left|\omega_{j+1}(\mu, w)-\omega_{j}(\mu, w)\right| \geq \frac{1}{c}-\frac{K}{j}>\frac{1}{2 c}
$$

if $j>K / 2 c$, uniformly in $\mu \in\left[\mu_{1}, \mu_{2}\right]$, $w \in B_{R}$. For $1 \leq j \leq K / 2 c$ the minimum

$$
m_{j}:=\min _{(\mu, w) \in\left[\mu_{1}, \mu_{2}\right] \times B_{R}}\left|\omega_{j+1}(\mu, w)-\omega_{j}(\mu, w)\right|
$$

is attained because $a \mapsto \lambda_{j}(a)$ is a compact function on $H^{1}$ by (68) and the compact embedding $H^{1}(0, \pi) \hookrightarrow L^{\infty}(0, \pi)$ (see also [20], Theorem 3 p. 31 and p.34). Each $m_{j}>0$ because all the eigenvalues $\lambda_{j}$ are simple.

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[^1]:    ${ }^{2} z^{*}$ denotes the complex conjugate of $z \in \mathbb{C}$.

[^2]:    ${ }^{3}$ In this case $\omega_{j}(\mu, \omega)=i \sqrt{\left|\lambda_{j}(\mu, \omega)\right|}$ is a purely imaginary complex number.

[^3]:    ${ }^{4}$ This means that equation (3) is non-autonomous indeed.

[^4]:    ${ }^{5}$ The operator norm is $\|U\|_{\sigma}:=\sup _{\|h\|_{\sigma} \leq 1}\|U h\|_{\sigma}$.

