# PROPERTIES OF SOME BREATHER SOLUTIONS OF A NONLOCAL DISCRETE NLS EQUATION 

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#### Abstract

We present results on breather solutions of a discrete nonlinear Schrödinger equation with a cubic Hartree-type nonlinearity that models laser light propagation in waveguide arrays that use a nematic liquid crystal substratum. A recent study of that model by Ben et al [BCMP15] showed that nonlocality leads to some novel properties such as the existence of orbitaly stable breathers with internal modes, and of shelf-like configurations with maxima at the interface. In this work we present rigorous results on these phenomena and consider some more general solutions. First, we study energy minimizing breathers, showing existence as well as symmetry and monotonicity properties. We also prove results on the spectrum of the linearization around one-peak breathers, solutions that are expected to coincide with minimizers in the regime of small linear intersite coupling. A second set of results concerns shelf-type breather solutions that may be thought of as limits of solutions examined in [BCMP15]. We show the existence of solutions with a non-monotonic front-like shape and justify computations of the essential spectrum of the linearization around these solutions in the local and nonlocal cases.


Key words. discrete nonlinear Schrödinger equations, nonlocal effects, localized solutions, breather solutions, linear stability

AMS subject classifications. 37K60, 35Q55, 47J10, 47J30, 78A60

## 1. Introduction

We present results on breather solutions of a nonlocal discrete nonlinear Schrödinger equation with a cubic Hartree-type nonlinearity introduced by Fratalocchi and Assanto [FA05] to study the propagation of laser light in waveguide arrays built from a nematic liquid crystal substratum. Further experimental work on this system is reported in [FABK04], [FABK05], [RAK13]. Breather solutions of discrete NLS equations have both theoretical and experimental importance and have been studied extensively, especially for the discrete NLS equation with power nonlinearities, see [K09] and references. Recently Ben et al [BCMP15] reported some interesting new properties of breather solutions of the nonlocal discrete NLS system of [FA05] in a 1-D lattice, specifically the presence of internal modes around orbitaly stable breathers, and the appearance of front-like solutions that attain their maximum at the interface. The present work studies some of these properties from a rigorous point of view, and includes new related results on breathers of both the local and nonlocal discrete cubic NLS equations.

We will show existence and stability results for two types of solutions of the Fratalocchi-Assanto [FA05] system. First, we study breather solutions that minimize the Hamiltonian (energy) over configurations with fixed $l^{2}$-norm (power). These are spatially localized solutions that could be in principle studied experimentally, see [FABK04], [FABK05], [RAK13] for some related localized solutions. The existence of the energy minimizers follows from a concentration-compactness argument that is similar to the one used in [W99], [M05], [P05], [S07]. We also show that the minimizer is even and decays monotonically away from its maximum. These two properties follow from discrete symmetrization and rearrangement arguments, see Theorem 2.2 for statements

[^0]and proofs. The notion of symmetrization for functions on the lattice uses the Fourier transform, while the monotonicity argument uses a discrete analogue of Riesz's rearrangement inequality, see [LL01].

We note that [BCMP15] shows the existence of spatially decaying breathers using different arguments, assuming that the linear intersite coupling parameter $\delta$, see (2.1) below, is small in absolute value. The present results assume instead that $\delta \gamma>0$, see (2.1), and thus extend the existence results of [BCMP15] to a larger parameter range. The $\delta \gamma>0$ case is also the one relevant to experiments.

The numerical results of [BCMP15] suggest that the linearization around the minimizing solutions of section 2 should have internal modes. Showing this in the general case seems difficult and we here present results for the linearization around one-peak breather solutions obtained for small linear intersite coupling $\delta$ via a continuation argument, see [BCMP15]. These solutions exist in the intersection of the parameter ranges considered in section 2 and in [BCMP15], and we argue that they coincide with the energy minimizers of section 2, see Proposition 4.5 for the precise statement. The presence of stable internal modes was seen explicitly in the limiting case on vanishing linear coupling $\delta$ in [BCMP15]. In section 4 we use perturbation arguments, especially the Krein signature, to show the persistence of at least a finite number of stable internal modes for $\delta>0$ sufficiently small, see Proposition 4.4. These results agree with the numerical observations of [BCMP15].

The second type of solutions we study are shelf-type breather solutions. These are solutions that decay in one direction and asymptote to a nonzero value in the other direction. These solutions have infinite $l^{2}-$ norm, and their existence is shown by an implicit function theorem that allows us to continue solutions of the $\delta=0$ equation, see [MA94], [PKF05a], [PP08], [BCMP15] for related results for solutions that decay at infinity. The argument applies to both the power and nonlocal nonlinearities, see section 3. In the nonlocal case we continue from a nontrivial configuration that exhibits a maximum at the interface between sites with vanishing and nonvanishing amplitude. This effect was seen originally in [BCMP15] for $l^{2}$ solutions, and we see here that the linearization can be still inverted explicitly.

We also consider the linearization around the shelf-type breathers of section 3 and present results on the essential spectrum of the corresponding operators. In the case of the power nonlinearity, the essential spectrum is obtained by studying the "spatial" dynamics of the piecewise constant coefficient linear map defined by the linear operator. This is a common approach for problems on the line, see e.g. [KP13], and we show discrete analogues of some of the basic results in subsection 4.2. In the nonlocal case we do not seem to have an immediate correspondence of the operator problem to spatial dynamics, but we can reproduce some of the calculations of the local case using the fact that the operator problem involves piecewise translation invariant (convolution) operators to obtain a weaker result for the essential spectrum, see Theorem 4.13 in subsection 4.3. A possible alternative approach may be provided by a study of the analytical properties of the resolvent, see [KKK06] for discrete Schrödinger operators, and [PS11] for the linearization around NLS breathers.

We note that the motivation for studying the shelf-type solutions is heuristic at present. First, the shelf-type breathers can be thought of as limits of a class of breathers that have finite support in the $\delta=0$ limit and were examined in [BCMP15]. These "finite shelf-type breather" solutions can be characterized by an integer $m>0$, the size of their support in the $\delta=0$ limit. The essential spectrum computed in section 4 is the apparent limit of a subset of the spectrum of the linearization around the finite shelf-
type breathers as $m$ increases, see details in section 4 . We also point out that there are interesting connections between shelf-type breathers of discrete NLS equations and static front solutions of discrete nonlinear diffusion equations. The equations for these solutions are similar, and the continuation methods of section 3 can be also used to show existence and monotonicity properties of static fronts, see [CCP15], [HM10]. In section 3 we remark on the discrete version of a recently proposed diffusion equation with nonlocal nonlinearity, see [FCET13]. Also, in discrete diffusion equations the static solutions can become "depinned", i.e. may start to move after some intersite coupling values, and the static problem is a starting point for analyzing these moving fronts [CB03]. In the discrete NLS we do not have a clear analogue of the depinning transition for shelf-type solutions, but a preliminary numerical study [P16] suggests that such a transition may be possible. This transition can be related to changes in the spectrum and bifurcations of the shelf-type breathers as we increase the linear coupling.

The Fratalocchi-Assanto model can be generalized to higher dimensions, see e.g. [W99], [SK05] for dimension dependent phenomena in the local system, however the experimental setup of [FA05], [FABK04], [FABK05], [RAK13] does not have a higher dimensional analogue. Also, comparisons of the model to experiments have used $\delta, \gamma, \kappa$ near unity, see [FA05], and thus the results of section 2 are directly relevant. The small $\delta$ results of sections 3,4 agree with the numerical results of [BCMP15], where $\delta$ was up to 0.5 . Their relevance to the experimental parameter range requires further numerical study.

The paper is organized as follows. In section 2 we study energy minimizing breathers and show existence, and symmetry and monotonicity properties. In section 3 we show the existence of shelf-like solutions for the local and nonlocal nonlinearities, in subsections 3.1, 3.2 respectively. In section 4 we study the spectra of the linearization around different breathers. In subsection 4.1 we study the linear stability of one-peak breather solutions in the small linear coupling regime and relate these solutions to the minimizers of section 2. In subsections 4.2, 4.3 we present results on the essential spectra of shelf-type solutions in the local, and nonlocal cases respectively.

## 2. Nonlocal discrete NLS and energy-minimizing breather solutions

We consider the one-dimensional discrete NLS equation

$$
\begin{equation*}
\dot{u}_{n}=\delta i\left(u_{n+1}+u_{n-1}-2 u_{n}\right)+2 \gamma \tanh \frac{\kappa}{2} i\left(\sum_{m \in \mathbb{Z}} e^{-\kappa|m-n|}\left|u_{m}\right|^{2}\right) u_{n} \quad n \in \mathbb{Z}, \tag{2.1}
\end{equation*}
$$

with $\delta, \gamma, \kappa$ real constants, $\kappa>0$.
Equation (2.1) can be written formally as the Hamiltonian system

$$
\begin{gather*}
\dot{u}_{n}=-i \frac{\partial H}{\partial u_{n}^{*}}, \quad n \in \mathbb{Z}, \quad \text { with }  \tag{2.2}\\
H=\delta \sum_{n \in \mathbb{Z}}\left|u_{n+1}-u_{n}\right|^{2}-\gamma \tanh \frac{\kappa}{2} \sum_{n \in \mathbb{Z}} \sum_{m \in \mathbb{Z}}\left|u_{m}\right|^{2} e^{-\kappa|m-n|}\left|u_{n}\right|^{2} . \tag{2.3}
\end{gather*}
$$

The quantity $P=\sum_{n \in \mathbb{Z}}\left|u_{n}\right|^{2}$ is formally a conserved quantity.
We consider solutions of (2.1) of the "breather" form $u_{n}=e^{-i \omega t} A_{n}$, with $\omega$ real, and $A: \mathbb{Z} \rightarrow \mathbb{C}$ independent of $t$. Such $A$ satisfies

$$
\begin{equation*}
-\omega A_{n}=\delta\left(A_{n+1}+A_{n-1}-2 A_{n}\right)+2 \gamma \tanh \frac{\kappa}{2}\left(\sum_{m \in \mathbb{Z}} e^{-\kappa|m-n|}\left|A_{m}\right|^{2}\right) A_{n}, \quad \forall n \in \mathbb{Z} \tag{2.4}
\end{equation*}
$$

Let $\|\cdot\|_{p}$ denote the $l^{p}$-norm for complex-valued functions on $\mathbb{Z}$. Let $Y=l^{2}(\mathbb{Z}, \mathbb{C})$ with the norm $\|\cdot\|=\|\cdot\|_{2}$.

Let $\left\{\hat{e}_{n}\right\}_{n \in \mathbb{Z}}$ denote the standard basis in $Y$.
Let $\lambda>0$ and define

$$
\begin{equation*}
\left.I_{\lambda}=\inf \{H(v): v \in Y), P(v)=\lambda\right\} . \tag{2.5}
\end{equation*}
$$

A configuration $A_{*} \in Y, P\left(A_{*}\right)=\lambda$, satisfying $H\left(A_{*}\right)=I_{\lambda}$ is referred to as a minimizer or ground state ( of $H$ at power $P=\lambda$ ). The set of minimizers of $H$ at power $\lambda$ is denoted as $\mathcal{M}_{\lambda}$. For $q \in \mathbb{Z}, \tau_{q} A$ denotes the translation of $A$ by $q$, i.e. $\left(\tau_{q} A\right)_{n}=A_{n-q}$. The operation of multiplying $A \in Y$ by a scalar $e^{i \phi}, \phi \in \mathbb{R}$, is referred to a global phase change. Translations and global phase changes leave $P$, and $H$ invariant.
Theorem 2.1. Let $\delta, \gamma>0$. Then for every $\lambda>0$, the set $\mathcal{M}_{\lambda}$ is nonempty, and $A \in \mathcal{M}_{\lambda}$ implies that $A$ satisfies (2.4).

Proof. We will only present a sketch of the proof. The existence part uses a discrete version of Lyon's concentration-compactness principle, following [S07], [P05], [M05], [W99]. The functional $H$ is easily seen to be bounded below, and has the form $H=\delta H_{2}-\gamma V$, with $H_{2}$, $V$ homogeneous quadratic and quartic functionals respectively. Both $H_{2}(u)$ and $V(u)$ are positive for any $u \neq 0$ in $X$. The discrete version of the concentration-compactness lemma, see [P05], sections 3, 4, states that there are three scenarios: splitting, vanishing, and convergence (up to translation), indicated by the number $\Gamma$ in [P05], (3.13). The goal is to show convergence, for any $\lambda>0$. To eliminate splitting, we use the subadditivity property (Lemma 4.3, [P05]), and a discrete analogue of Lemma 4.4 of [PM13]. These arguments are valid for all $\lambda>0$. To eliminate vanishing for a given $\lambda>0$ it suffices to find a test function $f \in X$, with $P(f)=\lambda$, and $H(f)<0$. We use functions

$$
\begin{equation*}
f=C g, \quad g \in X, \quad C=\sqrt{\lambda}\|g\|^{-1} \tag{2.6}
\end{equation*}
$$

Then $P(f)=\lambda$, and the condition $H(f)<0$ becomes

$$
\begin{equation*}
\frac{\delta\|D g\|^{2}\|g\|^{2}}{\gamma V(g)}<\lambda \tag{2.7}
\end{equation*}
$$

where $(D g)_{n}=g_{n+1}-g_{n}, n \in \mathbb{Z}$. We use test functions $(g)_{n}=e^{-\alpha|n|}, n \in \mathbb{Z}$, with $\alpha>0$ (as in [W99]) to see that

$$
\begin{equation*}
\|D g\|^{2}\|g\|^{2}=2+O(\alpha), \quad V(g) \geq\|g\|_{4}^{4}=\frac{1}{2 \alpha}+O(1), \quad \text { as } \quad \alpha \rightarrow 0^{+} \tag{2.8}
\end{equation*}
$$

Then, given any $\lambda>0$, the left hand side of (2.7) can be made smaller by taking $\alpha>0$ sufficiently small. It follows that a minimizing sequence will converge up to translations and global phase rotations, and that $\mathcal{M}_{\lambda}$ is nonempty for all $\lambda>0$.

We easily check that $H: Y \rightarrow \mathbb{R}$ is $C^{1}$, and that $H^{\prime}(A)-\omega P^{\prime}(A)=0$ implies that $A$ is a solution of (2.4).
Theorem 2.2. Let $\delta, \gamma>0$, and consider $\lambda>0$. Then $A \in \mathcal{M}_{\lambda}$ implies that there exist $\phi \in \mathbb{R}, q \in \mathbb{Z}$ such that $A=e^{i \phi} \tau_{q} \bar{A}$, with $\bar{A}$ real and positive, reflection symmetric and non-increasing, i.e. $\bar{A}_{n} \geq 0, \forall n \in \mathbb{Z}, \bar{A}_{-n}=\bar{A}_{n}, \forall n \in \mathbb{Z}$, and $\bar{A}_{a} \geq \bar{A}_{b}$, for all positive integers $a<b$.

We first define a symmetrization operation that transforms any $u \in Y=l^{2}(\mathbb{Z}, \mathbb{C})$ into a configuration $w \in Y$ that is real, non-negative and symmetric with respect to
the origin, and that moreover has the same $l^{2}$-norm and a smaller value of $H$. The symmetrization, and the proof of Lemma 2.3 below were suggested to us by Diego Rial.

For $f \in Y$, define $\tilde{f}: \mathbb{R} \rightarrow \mathbb{C}$ by

$$
\begin{equation*}
\tilde{f}(x)=\sum_{n \in \mathbb{Z}} f_{n} e^{i n x}, \quad x \in \mathbb{R} . \tag{2.9}
\end{equation*}
$$

$\tilde{f}$ is clearly $2 \pi$-periodic. By the $L^{2}$ theory for Fourier series, $\tilde{f}$ belongs to $L^{2}\left(S^{1}, \mathbb{C}\right)$, with $S^{1}$ the interval $[0,2 \pi]$ with $0,2 \pi$ identified, moreover (2.9) defines an isomorphism between $Y$ and $L^{2}\left(S^{1}, \mathbb{C}\right)$. The inverse is given by

$$
\begin{equation*}
f_{n}=\frac{1}{2 \pi} \int_{0}^{2 \pi} \tilde{f}(x) e^{-i n x} d x, \quad n \in \mathbb{Z} \tag{2.10}
\end{equation*}
$$

i.e. $\tilde{f}$ is the inverse Fourier transform of $f$. Moreover, we have Parseval's identity

$$
\begin{equation*}
\sum_{n \in \mathbb{Z}}\left|f_{n}\right|^{2}=\frac{1}{2 \pi} \int_{S^{1}}|\tilde{f}(x)|^{2} d x \tag{2.11}
\end{equation*}
$$

Consider $u \in Y$, with $\tilde{u}$ as in (2.9), and define $v, w$ as follows:

$$
\begin{equation*}
v_{n}=\frac{1}{2 \pi} \int_{S^{1}}|\tilde{u}(x)| e^{-i n x} d x, \quad n \in \mathbb{Z} \tag{2.12}
\end{equation*}
$$

and

$$
\begin{equation*}
w_{n}=\left|v_{n}\right|, \quad n \in \mathbb{Z} \tag{2.13}
\end{equation*}
$$

It follows from (2.11) that $v, w$ belong to $Y$, and that $\|w\|=\|v\|=\|u\|$.
Also, by (2.12) we have $v_{-n}=v_{n}^{*}, \forall n \in \mathbb{Z}$. Then (2.13) implies that

$$
\begin{equation*}
w_{-n}=w_{n} \geq 0, \quad \forall n \in \mathbb{Z} \tag{2.14}
\end{equation*}
$$

The configuration $w$ is therefore real, positive, and symmetric, with $\|w\|=\|u\|$. Lemma 2.3. Consider $u \in Y$, and $v, w \in Y$ defined as in (2.12), (2.13). Then, $H(w) \leq$ $H(v) \leq H(u)$.

Proof. First we examine the quartic part $V$ of $H$ in Fourier space. Letting $g_{n}=$ $e^{-\kappa|n|}, n \in \mathbb{Z}$, we compute

$$
\begin{equation*}
\tilde{g}(x)=\sum_{n \in \mathbb{Z}} e^{-\kappa|n|} e^{i n x}=\frac{1-e^{-2 \kappa}}{1-2 e^{-\kappa} \cos x+e^{-2 \kappa}} . \tag{2.15}
\end{equation*}
$$

We observe that $|\tilde{g}(x)|=\tilde{g}(x), \forall x \in \mathbb{R}$.
Also,

$$
\begin{align*}
V(u) & =\sum_{n \in \mathbb{Z}} \sum_{m \in \mathbb{Z}}\left|u_{n}\right|^{2} g_{n-m}\left|u_{m}\right|^{2}  \tag{2.16}\\
& =\frac{1}{8 \pi^{3}} \int_{0}^{2 \pi} \int_{0}^{2 \pi} \int_{0}^{2 \pi} \tilde{g}(x) \tilde{u}(y) \tilde{u}^{*}(y+x) \tilde{u}(z) \tilde{u}^{*}(z-x) d x d y d z  \tag{2.17}\\
& \leq \frac{1}{8 \pi^{3}} \int_{0}^{2 \pi} \int_{0}^{2 \pi} \int_{0}^{2 \pi} \tilde{g}(x)|\tilde{u}(y)|\left|\tilde{u}^{*}(y+x)\right||\tilde{u}(z)|\left|\tilde{u}^{*}(z-x)\right| d x d y d z \tag{2.18}
\end{align*}
$$

The definition of $v$ in (2.12) and the isomorphism between $Y$ and $L^{2}\left(S^{1}, \mathbb{C}\right)$ imply $|\tilde{u}(x)|=\tilde{v}(x)$, a.e. $x \in S^{1}$, so that the R.H.S. of (2.18) is $V(v)$. Therefore

$$
\begin{equation*}
V(u) \leq V(v) \tag{2.19}
\end{equation*}
$$

On the other hand, by the definition of $w$ in (2.13) we also have $V(v)=V(w)$, hence

$$
\begin{equation*}
V(u) \leq V(w) \tag{2.20}
\end{equation*}
$$

For the quadratic part $\mathrm{H}_{2}$ we calculate

$$
\begin{equation*}
H_{2}(u)=\frac{1}{\pi} \int_{0}^{2 \pi} \sin ^{2} \frac{x}{2}|\tilde{u}(x)|^{2} d x=H_{2}(v) \tag{2.21}
\end{equation*}
$$

since $|\tilde{u}(x)|=\tilde{v}(x)$, a.e. $x \in S^{1}$. Then

$$
\begin{equation*}
H_{2}(v)=\sum_{n \in \mathbb{R}}\left|v_{n+1}-v_{n}\right|^{2} \geq \sum_{n \in \mathbb{R}}| | v_{n+1}|-| v_{n} \|^{2}=H_{2}(w), \tag{2.22}
\end{equation*}
$$

using the definition of $w$ in (2.13). Combining (2.20), (2.22) we conclude that $H(u) \geq$ $H(w)$, as stated.

We also observe that by (2.12)

$$
\begin{equation*}
v_{0}=\frac{1}{2 \pi} \int_{0}^{2 \pi}|\tilde{u}(x)| d x \geq 0 ; \quad\left|v_{n}\right| \leq v_{0}, \quad \forall n \in \mathbb{Z} \backslash\{0\} \tag{2.23}
\end{equation*}
$$

Then (2.13) implies

$$
\begin{equation*}
w_{n} \leq w_{0}, \quad \forall n \in \mathbb{Z} \backslash\{0\} \tag{2.24}
\end{equation*}
$$

i.e. the maximum of $w$ is attained at the origin.

The second operation that decreases $H$ while keeping the $l^{2}$-norm fixed is that of replacing $u \in Y$ by its non-increasing rearrangement. To define this operation, let $\mathcal{S}$ be the set of all $f: \mathbb{Z} \rightarrow \mathbb{R}$ that satisfy (i) $f_{-n}=f_{n}, \forall n>0$, (ii) $f_{0}>0$, (iii) $f_{0} \geq f_{n} \geq 0$, $\forall n>0$, and (iv) $f_{n} \rightarrow 0$ as $|n| \rightarrow \infty$.

Note that $\mathcal{S}$ is the union of its disjoint subsets $\mathcal{S}^{+}, \mathcal{S}^{0}$, and $\mathcal{S}^{F}$, where $f \in \mathcal{S}^{+}$implies $f_{n}>0, \forall n \in \mathbb{Z}, f \in \mathcal{S}^{0}$ implies that there exist $n$ for which $f_{n}=0$ and that $f_{n} \neq 0$ for an infinite subset of $n \in \mathbb{Z}$, and $f \in \mathcal{S}^{F}$ implies that the set of $n$ for which $f_{n} \neq 0$ is finite.

For $f \in \mathcal{S}^{+} \cup \mathcal{S}^{0}$ as above, let $f^{j}$ the $j$-th largest value of $f$, starting from $j=1$, i.e. $f^{1}=f_{0}$, and $j \in \mathbb{Z}^{+}=\{1,2, \ldots\}$, for all $f \in \mathcal{S}$. For if $f \in \mathcal{S}^{F}$, then $f$ attains a finite number $\mathcal{T}(f)$ of values; then $f^{j}, j \leq \mathcal{T}(f)$ is the $j$-th largest value of $f$, starting from $j=1$, and we set $f^{j}=0$, for all $j>\mathcal{T}(f)$.

Given a subset $J$ of $\mathbb{Z}$, $\chi_{J}$ will denote the characteristic function of $J$. We also use the abbreviated notation $\chi_{f \geq a}$ to denote the characteristic function of the set of $n \in \mathbb{Z}$ satisfying $f_{n} \geq a$. Note that the support of $\chi_{f \geq f^{j}}$ is finite, for any $f \in \mathcal{S}$, and $j \geq 1$. In what follows $\chi_{J}(n)$ denotes the value of $\chi_{J}$ at site $n$.
Definition 2.4. Let $\mathcal{I}^{S}$ be the set of finite $I \subset \mathbb{Z}$ satisfying $0 \in I$, and $n \in I$ implies $-n \in I$. Let $m$ be such that $|I|=2 m+1$. Then $\bar{I}$ denotes the set of $2 m+1$ consecutive integers from $-m$ to $m$.

Let $f \in S^{+} \cup S^{F}$, fix $F^{j}>0$ and let $I=\left\{n \in \mathbb{Z}: f_{n} \geq f^{j}\right\}=\operatorname{supp} \chi_{f \geq f^{j}}$. Then $I \in \mathcal{I}^{S}$ and we let $\bar{\chi}_{f \geq f^{j}}=\chi_{\bar{I}}$.

For $f \in \mathcal{S}$ we define its symmetric non-increasing rearrangement as the function $\bar{f}: \mathbb{Z} \rightarrow \mathbb{R}^{+} \cup\{0\}$ defined by

$$
\begin{equation*}
\bar{f}=\sum_{j=1}^{\infty}\left(f^{j}-f^{j+1}\right) \bar{\chi}_{f \geq f^{j}} . \tag{2.25}
\end{equation*}
$$

We check that $\bar{f} \in \mathcal{S}^{+} \cup \mathcal{S}^{F}$, and that $\bar{f}_{m} \leq \bar{f}_{n}$ for all $m>n \geq 0$. We recall the following additional facts.
Proposition 2.5. If $f \in \mathcal{S}$ then

$$
\begin{equation*}
g_{n}=\sum_{j=1}^{\infty}\left(f^{j}-f^{j+1}\right) \chi_{f \geq f^{j}}(n), \quad n \in \mathbb{Z}, \tag{2.26}
\end{equation*}
$$

coincides with $f$, i.e. $g_{n}=f_{n}, \forall n \in \mathbb{Z}$. If also $f \in \mathcal{S} \cap l^{1}$, then

$$
\begin{equation*}
\sum_{n \in \mathbb{Z}} \bar{f}_{n}=\sum_{n \in \mathbb{Z}} f_{n} . \tag{2.27}
\end{equation*}
$$

Remark 2.6. The second conclusion of Proposition 2.5 also implies that $\|\bar{h}\|_{l^{p}}=\|h\|_{l^{p}}$, for all $h \in \mathcal{S} \cap l^{p}, p>0$. This follows by setting $f_{n}=\left(h_{n}\right)^{p}$.

Proof. First, suppose that $n$ is such that $f_{n}=f^{m}$, with $f^{j}>0$. Then $\chi_{f \geq f^{j}}(n)=1$, if $j \geq m$, and vanishes otherwise, and therefore

$$
\begin{equation*}
g_{n}=\sum_{j=m}^{\infty}\left(f^{j}-f^{j+1}\right)=\lim _{N \rightarrow \infty}\left(f^{m}-f^{N}\right)=f^{m}=f_{n} \tag{2.28}
\end{equation*}
$$

since $f^{j} \rightarrow 0$ as $j \rightarrow \infty$. In the case where $f_{n}=0$ for some $n$, then $g_{n}=0$ by (2.25), since $n$ is not in the support of any $\chi_{f \geq f^{j}}$ with $f^{j}>0$.

For the second statement, we first consider the case $f \in \mathcal{S}^{+}$. Let $\mathcal{N}\left(f^{j}\right), j \in \mathbb{Z}^{+}$be the set of $n$ satisfying $f_{n}=f^{j}$. The $\mathcal{N}\left(f^{j}\right)$ are finite, mutually disjoint, and their union is $\mathbb{Z}$. Moreover the set of $f^{j}, j \in \mathbb{Z}^{+}$is the set of values of $f$. Then

$$
\begin{equation*}
\sum_{n \in \mathbb{Z}} f_{n}=\sum_{j=1}^{\infty} f^{j}\left|\mathcal{N}\left(f^{j}\right)\right| \tag{2.29}
\end{equation*}
$$

moreover $\mathcal{N}\left(f^{j}\right)=\operatorname{supp} \chi_{f \geq f^{j}}$. Therefore (2.29) implies

$$
\begin{equation*}
\sum_{n \in \mathbb{Z}} f_{n}=\sum_{j=1}^{\infty} f^{j}\left|\chi_{f \geq f^{j}}\right|=\sum_{j=1}^{\infty} f^{j}\left|\bar{\chi}_{f \geq f^{j}}\right|, \tag{2.30}
\end{equation*}
$$

since the support of $\chi_{f \geq f^{j}}$ belongs to $\mathcal{I}^{S}$. We check that $\bar{f}^{j}=f^{j}$, and that $\bar{\chi}_{f \geq f^{j}}=$ $\chi_{\bar{f} \geq \bar{f}^{j}}, \forall j \in \mathbb{Z}^{+}$. Thus (2.30) implies

$$
\begin{equation*}
\sum_{n \in \mathbb{Z}} f_{n}=\sum_{j=1}^{\infty} \bar{f}^{j}\left|\chi_{\bar{f} \geq \bar{f}^{j}}\right|=\sum_{j=1}^{\infty} \bar{f}^{j} \mathcal{N}\left(\bar{f}^{j}\right)=\sum_{n \in \mathbb{Z}} \bar{f}_{n} . \tag{2.31}
\end{equation*}
$$

The case $f \in \mathcal{S}^{0}$ is similar. We note there that the set where $f_{n}=0$ is not included in the sum of (2.29), i.e. we only sum over the positive values of $f$. The set of positive
values of $f$ is precisely the set of the $f^{j}, j \in \mathbb{Z}^{+}$, and the remaining steps are identical since the sets $\mathcal{N}\left(f^{j}\right)$ satisfy the same properties. In the case $f \in \mathcal{S}^{F}$ we only sum over the values $f^{j}, j \leq \mathcal{T}(f)$ in (2.29). The remaining steps are identical as the $\mathcal{N}\left(f^{j}\right)$, $1 \leq j \leq \mathcal{T}(f)$, satisfy the same properties.

From the properties of the configuration $w \in Y$ obtained from $u \in Y$ via (2.12), (2.13), as above, it is clear that $w \in \mathcal{S} \cap l^{2}$, and we can apply the operation of nonincreasing rearrangement to obtain a new symmetric configuration $\bar{w}$ that is nonincreasing, and satisfies $\|\bar{w}\|=\|w\|$. To show that $H(\bar{w}) \leq H(w)$, we use Lemma 2.7 below on rearrangements of quadratic forms. The lemma is a special case of a discrete analogue of the Riesz rearrangement inequality, see e.g. [LL01].

Let $k$ be a positive integer and let $I, J$ be finite subsets of $\mathbb{Z}$, and let

$$
\begin{equation*}
Q(I, k, J)=\sum_{n \in \mathbb{Z}} \sum_{m \in \mathbb{Z}} \chi_{I}(n) \mathcal{H}_{k}(m-n) \chi_{J}(m), \tag{2.32}
\end{equation*}
$$

where $\mathcal{H}_{k}(n)=1$ if $|n| \leq k$, and vanishes otherwise.
Lemma 2.7. Let $I, J \in \mathcal{I}^{S}$, and consider the corresponding $\bar{I}, \bar{J}$, as in Definition 2.4. Assume also that $I \subseteq J$. Then for any positive integer $k$

$$
\begin{equation*}
Q(\bar{I}, k, \bar{J}) \geq Q(I, k, J) \tag{2.33}
\end{equation*}
$$

Given a pair $I, J$ the inequality is strict for at least one positive $k$, unless $\bar{I}=I, \bar{J}=J$.
Proof. Fix a positive integer $k$. Consider first the case $I=J$. Let $1 \leq j_{1}<j_{2}, \ldots,<j_{m}$ denote the positive sites of $J$. Then $|J|=2 m+1$. A gap $G$ of $J$ is a symmetric set of consecutive sites $j$ not in $J$ that satisfy either $j_{\lambda}<j<j_{\lambda+1}$ or $-j_{\lambda+1}<j<-j_{\lambda}$ for a pair of $0<j_{\lambda}<j_{\lambda+1}$ in $J$. Note that the definition implies that $g=j_{\lambda+1}-j_{\lambda}>1$, and that $G$ consists of $2 g>0$ sites. A gap is specified by the pair of sites $j_{\lambda}, j_{\lambda+1}$ above, referred to below as its (positive) edges.

If $J=\bar{J}$, then $J$ has no gaps. The idea is to show that deleting gaps in $I$ and $J$ simultaneously increases $Q$. We then generalize to the case $I \subset J$.

Given a gap $G$ specified by the sites $j_{\lambda}, j_{\lambda+1}$ of $J$ as above, let $J_{\operatorname{sub}} G \in \mathcal{I}^{S}$ be the set whose positive elements are all $j \in J$ that satisfy $j \leq j_{\lambda}$, and all $j-g, j \in J$ that satisfy $j \geq j_{\lambda}+g+1$. It is easy to show that $Q(J \operatorname{sub} G, k, J \operatorname{sub} G) \geq Q(J, k, J)$. Iterating the elimination of gaps we obtain $Q(\bar{J}, k, \bar{J}) \geq Q(J, k, J)$.

In the case $I \subseteq J$, the definition of gaps is the same. We first eliminate all gaps of $J$ whose positive sites are to the right of the largest site of $I$. The resulting set is denoted as $J_{1}$, and we easily argue that $Q\left(I, k, J_{1}\right) \geq Q(I, k, J)$. Also $I \subseteq J_{1}$, moreover the remaining gaps of $J_{1}$ are also gaps of $I$, otherwise there would be sites of $I$ not in $J_{1}$. Eliminating all the gaps of $J_{1}$ from both $I$ and $J_{1}$ we obtain new sets $I_{2}$, and $J_{2}$ respectively. We have $J_{2}=\bar{J}$, and we easily show that $Q\left(I_{2}, k, \bar{J}\right) \geq Q\left(I_{2}, k, J_{1}\right)$. Also $I_{2} \subseteq \bar{J}$. If $I_{2} \neq \bar{J}$, then $I_{2}$ still has gaps, and we eliminate them to obtain $\bar{I}$. We again can show that $Q(\bar{I}, k, \bar{J}) \geq Q\left(I_{2}, k, \bar{J}\right)$. The above hold for any positive integer $k$.

It is easy to show in each of the above steps, that cancellation of any gap would give strict inequality for all values of $k \geq 1$. For instance, the relation $I \subseteq J$ at all stages implies that the elimination of the common points in the gap will increase for each $i$ in the smaller set the number of $j$ in the larger set satisfying $|i-j| \leq k$. The increase will be strict for the points at the edges of the gap, for all $k \geq 1$. $\square$
Lemma 2.8. Consider the configuration $w \in \mathcal{S} \cap Y$, and its non-increasing symmetrization $\bar{w} \in Y$, as above. Then $H(\bar{w}) \leq H(w)$. The inequality is strict, unless $w=\bar{w}$.

Proof. We first consider the quartic part $V$, using the functions $f_{n}=w_{n}^{2}, \bar{f}_{n}=\bar{w}_{n}^{2}$, $\forall n \in \mathbb{Z}$, i.e. $f, \bar{f} \in \mathcal{S} \cap l^{1}$. Also let $g(n)=e^{-\kappa|n|}, \forall n \in \mathbb{Z}$. We can use Proposition 2.5 and
the definition $\mathcal{H}_{k}$ after (2.32) to write $g=\sum_{j=1}^{\infty}\left(g^{j}-g^{j+1}\right) \mathcal{H}_{j}$. Then,

$$
\begin{align*}
V(u) & =\sum_{n \in \mathbb{Z}} \sum_{m \in \mathbb{Z}} f_{n} g(n-m) f_{m}  \tag{2.34}\\
& =\sum_{n \in \mathbb{Z}} \sum_{m \in \mathbb{Z}}\left[\sum_{j_{1}=1}^{\infty} a_{j_{1}} \chi_{f \geq f^{j_{1}}}(n)\right]\left[\sum_{j=1}^{\infty} b_{j} \mathcal{H}_{j}(m-n)\right]\left[\sum_{j_{2}=1}^{\infty} a_{j_{2}} \chi_{f \geq f^{j_{2}}}(m)\right], \tag{2.35}
\end{align*}
$$

with $a_{j}=f^{j}-f^{j+1}, b_{j}=g^{j}-g^{j+1}$.
The order of summations can be changed using the $l^{2}$-summability of $w$, and the fact that convolution with $g$ defines a bounded operator in $Y$ (we omit therefore the steps of defining all quantities for finite sums and taking limits).

We then have

$$
\begin{align*}
V(u) & =\sum_{j_{1}=1}^{\infty} \sum_{j=1}^{\infty} \sum_{j_{2}=1}^{\infty} a_{j_{1}} b_{j} a_{j_{2}}\left[\sum_{n \in \mathbb{Z}} \sum_{m \in \mathbb{Z}} \chi_{f \geq f_{1}^{j}}(n) \mathcal{H}_{j}(m-n) \chi_{f \geq f_{2}^{j}}(m)\right]  \tag{2.36}\\
& =\sum_{j_{1}=1}^{\infty} \sum_{j=1}^{\infty} \sum_{j_{2}=1}^{\infty} a_{j_{1}} b_{j} a_{j_{2}} Q\left(\operatorname{supp} \chi_{f \geq f^{j_{1}}, j, \operatorname{supp}} \chi_{f \geq f^{j_{2}}}\right) . \tag{2.37}
\end{align*}
$$

Since $\operatorname{supp} \chi_{f \geq f^{j}} \subseteq \operatorname{supp} \chi_{f \geq f^{j^{\prime}}}$ if $f^{j} \geq f^{j^{\prime}}$, we can apply Lemma (2.7) for the supports of the $\chi_{f \geq f^{j}}$ to obtain from (2.37)

$$
\begin{align*}
V(f) & \geq \sum_{j_{1}=1}^{\infty} \sum_{j=1}^{\infty} \sum_{j_{2}=1}^{\infty} a_{j_{1}} b_{j} a_{j_{2}} Q\left(\operatorname{supp} \bar{\chi}_{f \geq f^{j_{1}}}, j, \operatorname{supp} \bar{\chi}_{f \geq f^{j_{2}}}\right)  \tag{2.38}\\
& =V(\bar{f}) . \tag{2.39}
\end{align*}
$$

The inequality is strict unless

$$
\begin{equation*}
Q\left(\operatorname{supp} \chi_{\left.f \geq f^{j_{1}}, j, \operatorname{supp} \chi_{f \geq f^{j_{2}}}\right)=Q\left(\operatorname{supp} \bar{\chi}_{f \geq f^{j_{1}}}, j, \operatorname{supp} \bar{\chi}_{f \geq f^{j_{2}}}\right), ., ~}\right. \tag{2.40}
\end{equation*}
$$

for all $j_{1}, j_{2}$. In that case we see that $\bar{w}=w$.
The quadratic part $H_{2}$ can be written as

$$
\begin{equation*}
H_{2}(w)=\sum_{n \in \mathbb{Z}}\left(w_{n+1}^{2}+w_{n}^{2}-2 w_{n} w_{n+1}\right)=2\|w\|^{2}-2 \sum_{n \in \mathbb{Z}} w_{n} w_{n+1}, \tag{2.41}
\end{equation*}
$$

and we define $H_{2}^{+}$by

$$
\begin{equation*}
H_{2}^{+}(w)=\sum_{n \in \mathbb{Z}}\left(w_{n+1}^{2}+w_{n}^{2}+2 w_{n} w_{n+1}\right)=2\|w\|^{2}+2 \sum_{n \in \mathbb{Z}} w_{n} w_{n+1} . \tag{2.42}
\end{equation*}
$$

To show that $H_{2}(\bar{w}) \leq H_{2}(w)$ it then suffices to show that $H_{2}^{+}(\bar{w}) \geq H_{2}^{+}(w)$. We write

$$
\begin{equation*}
H_{2}^{+}(w)=\sum_{n \in \mathbb{Z}} \sum_{m \in \mathbb{Z}} w_{n} M(n-m) w_{m}, \tag{2.43}
\end{equation*}
$$

with $M(0)=2, M( \pm 1)=1$, and $M(n)=0$ for $n \in \mathbb{Z} \backslash\{-1,0,1\} . M$ is symmetric and decreasing in $\mathbb{Z}^{+}$, and we can then follow the argument used for $V$ above to see that $H_{2}^{+}(\bar{w}) \geq H_{2}^{+}(w)$, as required. It follows immediately by $H(\bar{w})=H_{2}(\bar{w})-V(\bar{w}) \leq$ $H_{2}(w)-V(w)=H(w)$, with strict inequality if $\bar{w} \neq w$, as stated.

Remark 2.9. In the case $\delta, \gamma<0$ we have a similar statement for the maximizer of $H$, i.e. we apply the theorem to the minimizer of $-H$.

Remark 2.10. For $\delta>0, \gamma<0$ and any $\lambda>0$ we have $I_{\lambda} 0$. We can see that considering the functions $g^{N}$, $N \in \mathbb{Z}^{+}$, where $g_{n}^{N}=\sqrt{\lambda} / \sqrt{(2 N+1)}$ if $-N \leq n \leq N$, and vanishes elsewhere. It follows the infimum is not attained and that we can not use mimimization to find critical points. In the case of $\delta$ sufficiently, small critical points were obtained different arguments, see [BCMP15].
REMARK 2.11. The symmetrization and rearrangement results above also apply to energy minimizing breathers of the cubic DNLS in $\mathbb{Z}$, where we minimize $H_{q}=\delta H_{2}$ $\gamma H_{4}, H_{4}(u)=\|u\|_{l^{4}}^{4}$, over configurations $u \in Y$ with $P(u)=\lambda>0$. The existence of a minimizer for arbitrary $\lambda>0$ is shown in [W99]. Lemma 2.3 also holds for $H_{4}$, using the $k \rightarrow \infty$ limit $\tilde{g}(x)=1, \forall x$, in (2.15). Lemma 2.8 clearly applies to $H_{q}$ as well.

## 3. Infinite shelf-type breathers

In this section we show the existence of infinite shelf-type breather solutions of (2.4) with $|\delta| \neq 0$ and small. We consider two cases depending on the interaction term in (2.1): local case and nonlocal case, corresponding to " $\kappa=\infty$ ", and $\kappa>0$ respectively. In Theorem 3.2 we show the continuation of shelf-type breathers of (2.4) with $\delta=0$ in the local case. Theorem 3.3 concerns the continuation of shelf-type solutions for arbitrary $\kappa>0$.

### 3.1. Local Case

Let $\omega \in \mathbb{R}$. For $A \in l^{\infty}(\mathbb{Z}, \mathbb{R}), \delta \in \mathbb{R}$ and $n \in \mathbb{Z}$ define $F_{n}: l^{\infty}(\mathbb{Z}, \mathbb{R}) \times \mathbb{R} \longrightarrow \mathbb{R}$ by

$$
\begin{equation*}
F_{n}(A, \delta)=\delta(\Delta A)_{n}+2 \gamma A_{n}^{3}+\omega A_{n} \tag{3.1}
\end{equation*}
$$

where $(\Delta A)_{n}=A_{n-1}-2 A_{n}+A_{n+1}$.
Let $\alpha>0$ and define $S_{0}=\left\{n_{0}+1, n_{0}+2, n_{0}+3, \ldots\right\} \subset \mathbb{Z}, S_{A}=\mathbb{Z} \backslash S_{0}$. We consider a nontrivial solution $A=\widetilde{A}$ of $F_{n}(A, 0)=0, \forall n \in \mathbb{Z}$, obtained by requiring $A_{n} \neq 0$ for $n \in S_{A}$, and $A_{n}=0$ for $n \in S_{0}$. As $\delta=0$, the equation $F_{n}=0$ becomes

$$
\begin{equation*}
-2 \gamma\left|A_{n}\right|^{2}=\omega, \quad \forall n \in S_{A} \tag{3.2}
\end{equation*}
$$

Then $\gamma$ and $\omega$ must have opposite signs and we consider the solution

$$
\widetilde{A}_{n}=\left\{\begin{array}{l}
0, \text { if } n \in S_{0},  \tag{3.3}\\
\alpha, \text { if } n \in S_{A},
\end{array}\right.
$$

with $\alpha=\sqrt{-\frac{\omega}{2 \gamma}}$.
Remark 3.1. We can obtain more nontrivial solutions by taking $A_{n}<0$ for all $n \in S_{A}$. Moreover, we can choose $A_{n} \neq 0$ for $n \in S_{A}:=S_{+} \cup S_{-}$, with $A_{n}>0$ for $n \in S_{+}$and $A_{n}<$ 0 for $n \in S_{-}$provided $S_{+}$is finite or $S_{-}$is finite. The arguments below apply to these solutions as well.

Let $X=l^{2}(\mathbb{Z}, \mathbb{R})$ with the norm $\|\cdot\|=\|\cdot\|_{X}$ as in Section 2. Let $\|\cdot\|_{Z, Y}$ denote the operator norm of a linear operator from $Z$ to $Y$, both Banach spaces.

We define $G: X \times \mathbb{R} \rightarrow X$ as

$$
\begin{equation*}
G(B, \delta)=\left\{G_{n}(B, \delta)\right\}_{n \in \mathbb{Z}}=F_{n}(B+\widetilde{A}, \delta), \tag{3.4}
\end{equation*}
$$

with $F_{n}$ given by (3.1) It can be checked that $G$ is well-defined, i.e. if $B \in X, \delta \in \mathbb{R}$, then $G(B, \delta) \in X$.

Theorem 3.2. Let $G$ be as in (3.4) and consider the solution $\left(x_{0}, y_{0}\right)=(0,0)$ of $G(B, \delta)=0$. Then there exists $\epsilon_{0}>0$ and a unique continuous function $B:\left(-\epsilon_{0}, \epsilon_{0}\right) \rightarrow X$ satisfying $B(0)=0$ and $G(B(\delta), \delta)=0$ for all $\delta \in\left(-\epsilon_{0}, \epsilon_{0}\right)$.

The proof uses the implicit function theorem and is similar to that of Theorem 3.3 below (with simpler computations) and is omitted.

### 3.2. Nonlocal Case

We now show existence results for arbitrary $\kappa>0$ and $|\delta|$ small enough. Let $\omega \in \mathbb{R}$ and $\kappa>0$. For $A \in l^{\infty}(\mathbb{Z}, \mathbb{R})$ and $\delta \in \mathbb{R}$ define $F_{n}, n \in \mathbb{Z}$, by

$$
\begin{equation*}
F_{n}(A, \delta)=\delta\left(\Delta A_{n}\right)+2 \gamma \tanh \frac{\kappa}{2}\left(\sum_{m \in \mathbb{Z}} e^{-\kappa|m-n|}\left|A_{m}\right|^{2}\right) A_{n}+\omega A_{n} \tag{3.5}
\end{equation*}
$$

We consider, as before, a nontrivial solution $A=\widetilde{A}$ of $F_{n}(A, 0)=0, \forall n \in \mathbb{Z}$, determined by $\alpha>0$, and sets $S_{0}=\left\{n_{0}+1, n_{0}+2, n_{0}+3, \ldots\right\} \subset \mathbb{Z}, S_{A}=\mathbb{Z} \backslash S_{0}$, for which $A_{n}>0$ if $n \in S_{A}$, and $A_{n}=0$ if $n \in S_{0}$. Evaluating $\delta=0$ in (3.5), the equation $F_{n}=0$ becomes

$$
\begin{equation*}
-2 \gamma \tanh \frac{\kappa}{2} \sum_{m \in S_{A}} e^{-\kappa|m-n|}\left|A_{m}\right|^{2}=\omega, \quad \forall n \in S_{A} . \tag{3.6}
\end{equation*}
$$

We also write (3.6) as

$$
\begin{equation*}
M J=\left(-2 \gamma \tanh \frac{\kappa}{2}\right)^{-1} \omega \varepsilon, \tag{3.7}
\end{equation*}
$$

where $J_{m}=A_{m}^{2}, \varepsilon=[\ldots, 1,1,1]^{T} \in \mathbb{R}^{\left|S_{A}\right|}$ and $M$ is defined implicitly by the previous equation. Letting $\rho=e^{-\kappa}$ we have $M$ and its inverse $M^{-1}$ explicitly as

$$
M=\left[\begin{array}{lllll}
\cdots & & & \cdots  \tag{3.8}\\
& 1 & \rho & \rho^{2} & \rho^{3} \\
\rho & 1 & \rho & \rho^{2} \\
\rho^{2} & \rho & 1 & \rho \\
\cdots & \rho^{3} & \rho^{2} & \rho & 1
\end{array}\right], \quad M^{-1}=\frac{1}{1-\rho^{2}}\left[\begin{array}{ccccc}
\cdots & & & & \cdots \\
1+\rho^{2} & -\rho & 0 & 0 \\
-\rho & 1+\rho^{2} & -\rho & 0 \\
& 0 & -\rho & 1+\rho^{2} & -\rho \\
\cdots & 0 & 0 & -\rho & 1
\end{array}\right] .
$$

The calculation of $M^{-1}$ used the fact that $M$ is Toeplitz. It is clear from (3.8) that $M^{-1}$ is also bounded in $X$. Hence the vector $J$ is computed explicitly as

$$
J=\left(-2 \gamma \tanh \frac{\kappa}{2}\right)^{-1} \omega M^{-1} \varepsilon, \quad \text { where } \quad M^{-1} \varepsilon=\frac{1-\rho}{1+\rho}\left[\ldots, 1,1,1, \frac{1}{1-\rho}\right]^{T}
$$

Letting $\alpha^{2}=\left(-2 \gamma \tanh \frac{\kappa}{2}\right)^{-1} \omega \frac{1-\rho}{1+\rho}$, the solution $\widetilde{A}$ determined by $\alpha>0$ is

$$
\widetilde{A}_{n}=\left\{\begin{array}{ccc}
\alpha, & \text { if } & n \in S_{A},  \tag{3.9}\\
\frac{\alpha}{\sqrt{1-\rho}}, & , \text { if } n \in n_{0} \\
0, & \text { if } & n=n_{0}, \\
0, S_{0}
\end{array}\right.
$$

In a similar way we can find nontrivial solutions such that $A_{n}<0$ for all $n \in S_{A}$. Moreover, as in Remark 3.1, we can find nontrivial solutions such that $A_{n} \neq 0$ for $n \in$ $S_{A}:=S_{+} \cup S_{-}$, with $A_{n}>0$ for $n \in S_{+}$and $A_{n}<0$ for $n \in S_{-}$provided $S_{+}$is finite or $S_{-}$ is finite. The arguments below are valid for these solutions as well.

Then we define $G: X \times \mathbb{R} \rightarrow X$ as

$$
\begin{equation*}
G(B, \delta)=\left\{G_{n}(B, \delta)\right\}_{n \in \mathbb{Z}}=\left\{F_{n}(B+\widetilde{A}, \delta)\right\}_{n \in \mathbb{Z}} \tag{3.10}
\end{equation*}
$$

The fact that $G$ is well-defined follows from the arguments in the proof of Lemma 3.7 below.
Theorem 3.3. Let $G$ be as in (3.10) and fix $\omega, \kappa>0$. Let $n_{0}$ be a positive integer, and consider a nontrivial solution $(\widetilde{A}, 0)$ of $G(B, \delta)=0$ with $S_{A}=\left\{\ldots, n_{0}-2, n_{0}-1, n_{0}\right\}$. Then there exists $\epsilon_{0}>0$ and a unique continuous function $B:\left(-\epsilon_{0}, \epsilon_{0}\right) \rightarrow X$ satisfying $B(0)=0$ and $G(B(\delta), \delta)=0$ for all $\delta \in\left(-\epsilon_{0}, \epsilon_{0}\right)$.

Proof. As in the previous theorem, we apply the implicit function theorem around a solution ( $\widetilde{A}, 0)$ of $G(B, \delta)=0$. We present here the (Fréchet) derivative $D_{1} G$ and we show that it is a linear isomorphism of $X$. The continuity of $G$ and $D_{1} G$ at $(0,0)$ is proved in Lemmas 3.7 and 3.8 below. We have

$$
D_{1} G(0,0)=\left(\begin{array}{cc}
M_{1} & 0  \tag{3.11}\\
0 & M_{2}
\end{array}\right)
$$

with

$$
M_{1}=2 c \alpha^{2}\left[\begin{array}{ccccc}
\cdots & & & & \cdots \\
& 1 & \rho & \rho^{2} & \frac{\rho^{3}}{\sqrt{1-\rho}} \\
& \rho & 1 & \rho & \frac{\rho^{2}}{\sqrt{1-\rho}} \\
& \rho^{2} & \rho & 1 & \frac{\rho}{\sqrt{1-\rho}} \\
\cdots & \frac{\rho^{3}}{\sqrt{1-\rho}} & \frac{\rho^{2}}{\sqrt{1-\rho}} & \frac{\rho}{\sqrt{1-\rho}} & \frac{1}{1-\rho}
\end{array}\right], \quad M_{2}=\omega\left[\begin{array}{cccc}
1-\rho & 0 & 0 & \cdots \\
0 & 1-\rho^{2} & 0 & \cdots \\
0 & 0 & 1-\rho^{3} & \\
\cdots & & & \cdots
\end{array}\right],
$$

and $\rho=e^{-\kappa}, c=2 \gamma \tanh \frac{\kappa}{2}$.
As $M_{2}$ is an infinite diagonal matrix with $1-\rho \leq M_{2}(n, n)<1, \forall n \in \mathbb{Z}$, to prove that $D_{1} G$ is a linear isomorphism in $X$ it is enough to show that $M_{1}$ is invertible. We see that

$$
M_{1}^{-1}=\frac{1}{\left(c \alpha^{2}\right)\left(1-\rho^{2}\right)}\left[\begin{array}{ccccc}
\cdots & & & \cdots \\
1+\rho^{2} & -\rho & 0 & 0 \\
-\rho & 1+\rho^{2} & -\rho & 0 \\
& 0 & -\rho & 1+\rho^{2} & \frac{-\rho}{\sqrt{1-\rho}} \\
\ldots & 0 & 0 & \frac{-\rho}{\sqrt{1-\rho}} & 1-\rho
\end{array}\right],
$$

with $M_{1}^{-1}$ clearly bounded in $X$.
Remark 3.4. Observe that the nonlocal case approaches the local case as $\kappa \rightarrow \infty$, and $\rho=e^{-\kappa} \rightarrow 0$.
Remark 3.5. We note that [FCET13] propose a discrete nonlinear diffusion equation whose static front solutions satisfy (3.6) with the term $(1-\omega) A_{n}+A_{n}^{2}$ added to the right hand side (one may assume $A_{n}>0$ for all $n \in S_{A}$ ). The numerical results of [FCET13] suggest that this variant of (3.6) also has a solution with a maximum at the interface.

To prove the continuity of $G$ and $D_{1} G$ at $(0,0)$ we will need to use the following discrete version of Young's inequality:
Theorem 3.6. (Young's inequality, [LL01], ch. 4) Let p, $q, r \geq 1, \frac{1}{p}+\frac{1}{q}>1$ and $\frac{1}{p}+\frac{1}{q}=1+\frac{1}{r}$. Then $u \in l^{p}$, and $v \in l^{q}$ imply $u * v \in l^{r}$ and

$$
\begin{equation*}
\|u * v\|_{r} \leq\|u\|_{p}\|v\|_{q} \tag{3.12}
\end{equation*}
$$

The inequality follows from the statement for functions on the line, see e.g. [LL01], ch. 4.
Lemma 3.7. Let $G$ be as in (3.10) and fix $\omega$, and $\kappa>0$. Then $G$ is continuous at $(0,0) \in X \times \mathbb{R}$.

Proof. Let $(B, \delta) \in X \times \mathbb{R}$, then

$$
\begin{align*}
\|G(B, \delta)-G(0,0)\|_{X}^{2} & =\sum_{n \in \mathbb{Z}}\left|F_{n}(B+\tilde{A}, \delta)-F_{n}(\tilde{A}, 0)\right|^{2} \\
& =\sum_{n \in \mathbb{Z}}\left|\delta \Delta\left(B_{n}+\widetilde{A}_{n}\right)+c_{\gamma, \kappa}\left(B_{n}+\widetilde{A}_{n}\right) \sum_{m \in \mathbb{Z}} e^{-\kappa|m-n|}\right| B_{m}+\left.\widetilde{A}_{m}\right|^{2} \\
& +\omega\left(B_{n}+\widetilde{A}_{n}\right)-c_{\gamma, \kappa}\left(\sum_{m \in \mathbb{Z}} e^{-\kappa|m-n|} \widetilde{A}_{m}^{2}\right) \widetilde{A}_{n}-\left.\omega \widetilde{A}_{n}\right|^{2} \tag{3.13}
\end{align*}
$$

with $c_{\gamma, \kappa}=2 \gamma \tanh \frac{\kappa}{2}$, and

$$
\begin{equation*}
C_{n}=\sum_{m \in \mathbb{Z}} e^{-\kappa|m-n|}\left(\left|B_{m}+\widetilde{A}_{m}\right|^{2}\left(B_{n}+\widetilde{A}_{n}\right)-\widetilde{A}_{m}^{2} \widetilde{A}_{n}\right), \quad \forall n \in \mathbb{Z} \tag{3.14}
\end{equation*}
$$

Then

$$
\begin{equation*}
\|G(B, \delta)-G(0,0)\|_{X}^{2} \leq 6\left(|\delta|^{2}\|\Delta\|_{X, X}^{2}\|B\|_{X}^{2}+|\omega|\|B\|_{X}^{2}+2|\gamma| \tanh \frac{\kappa}{2}\|C\|_{X}^{2}\right) \cdot( \tag{3.15}
\end{equation*}
$$

The first and second term in the last inequality vanish as $\delta \rightarrow 0$ and $B \rightarrow 0 \in X$. It is then enough to show that $\|C\|_{X}$ vanishes as $\delta \rightarrow 0$ and $B \rightarrow 0 \in X$. We have

$$
\begin{align*}
\|C\|_{X}^{2} & =\sum_{n \in \mathbb{Z}}\left|\sum_{m \in \mathbb{Z}} \rho^{|m-n|}\left(\left|B_{m}+\widetilde{A}_{m}\right|^{2}\left(B_{n}+\widetilde{A}_{n}\right)-\widetilde{A}_{m}^{2} \widetilde{A}_{n}\right)\right|^{2}  \tag{3.16}\\
& =\sum_{n \in \mathbb{Z}}\left|\sum_{m \in \mathbb{Z}} \rho^{|m-n|}\right| B_{m}+\left.\widetilde{A}_{m}\right|^{2} B_{n}+\left.\sum_{m \in \mathbb{Z}} \rho^{|m-n|}\left(\left|B_{m}+\widetilde{A}_{m}\right|^{2}-\widetilde{A}_{m}^{2}\right) \widetilde{A}_{n}\right|^{2}  \tag{3.17}\\
& \leq \sum_{n \in \mathbb{Z}}\left|\sum_{m \in \mathbb{Z}} \rho^{|m-n|}\right| B_{m}+\left.\widetilde{A}_{m}\right|^{2} B_{n}+\left.\sum_{m \in \mathbb{Z}} \rho^{|m-n|} B_{m}\left(B_{m}+2 \widetilde{A}_{m}\right) \widetilde{A}_{n}\right|^{2}  \tag{3.18}\\
& \leq \Gamma\left(| | B \|\left|,||A||_{\infty}\right) \sum_{n \in \mathbb{Z}}\left|\sum_{m \in \mathbb{Z}} \rho^{|m-n|} B_{n}+\sum_{m \in \mathbb{Z}} \rho^{|m-n|} B_{m}\right|^{2}\right. \tag{3.19}
\end{align*}
$$

with

$$
\begin{equation*}
\Gamma\left(\|B\|,\|A\|_{\infty}\right)=4\left(\|B\|^{2}+\|A\|_{\infty}^{2}\right)^{2}+\|A\|_{\infty}^{2}\left(2\|B\|^{2}+4\|A\|_{\infty}^{2}\right) \tag{3.20}
\end{equation*}
$$

Therefore

$$
\begin{align*}
\|C\|_{X}^{2} & \leq 2 \Gamma\left(\|B\|,\|A\|_{\infty}\right)\left(\sum_{n \in \mathbb{Z}}\left|B_{n}\right|^{2}\left|\sum_{m \in \mathbb{Z}} \rho^{|m-n|}\right|^{2}+\sum_{n \in S_{A}}\left|\sum_{m \in \mathbb{Z}} \rho^{|m-n|} B_{m}\right|^{2}\right)  \tag{3.21}\\
& \leq 2 \Gamma\left(\|B\|,\|A\|_{\infty}\right)\left(\|B\|_{X}+\left\|\rho^{|\cdot|} * B\right\|_{X}\right)  \tag{3.22}\\
& \leq 2 \Gamma\left(\|B\|,\|A\|_{\infty}\right)\left(\|B\|_{X}+\left\|\rho^{\mid \cdot \cdot}\right\|_{1}\|B\|_{X}\right) . \tag{3.23}
\end{align*}
$$

The last inequality follows from Young's inequality (3.12) with $u=\rho^{|\cdot|}, v=B, r=2$, $p=1$ and $q=2$. By (3.20), (3.23) $\|C\|_{X}$ tends to zero if $\|B\|_{X}$ tends to zero.
Lemma 3.8. Fix $\omega$ and $\kappa>0$. and consider the function $G: X \times \mathbb{R} \rightarrow X$ defined in (3.10). Then $D_{1} G$ is continuous at $(\widetilde{A}, 0)$.

Proof. It is enough to show that there exists $\beta$ such that

$$
\begin{equation*}
\left\|\left[D_{1} G(B, \delta)-D_{1} G(0,0)\right] v\right\| \leq \beta\|v\|, \quad \forall v \in X, \tag{3.24}
\end{equation*}
$$

and $\beta \rightarrow 0$ as $(B, \delta) \rightarrow 0 \in X \times \mathbb{R}$. We have already determined $D_{1} G(0,0)$ in the proof of 3.2. We now calculate $D_{1} G(B, \delta)$ with $B \neq 0$, and $\delta \neq 0$, which (we check) is given by its partial derivatives. For $n \in \mathbb{Z}$ we have

$$
\begin{align*}
& \frac{\partial G_{n}}{\partial B_{n}}(B, \delta)=-2 \delta+2 \gamma \tanh \frac{\kappa}{2}\left(\sum_{m \in \mathbb{Z}} \rho^{-|m-n|}\left|B_{m}+\widetilde{A}_{m}\right|^{2}+2\left(B_{n}+\widetilde{A}_{n}\right)^{2}\right)+\omega,  \tag{3.25}\\
& \frac{\partial G_{n}}{\partial B_{m}}(B, \delta)=\delta+4 \gamma \tanh \frac{\kappa}{2}\left(\rho^{|m-n|}\left|B_{m}+\widetilde{A}_{m}\right|\right)\left(B_{n}+\widetilde{A}_{n}\right), \quad m \in\{(n-1),(n+1)\}, \tag{3.26}
\end{align*}
$$

$$
\begin{equation*}
\frac{\partial G_{n}}{\partial B_{m}}(B, \delta)=4 \gamma \tanh \frac{\kappa}{2}\left(\rho^{|m-n|}\left|B_{m}+\widetilde{A}_{m}\right|\right)\left(B_{n}+\widetilde{A}_{n}\right), \quad m \in \mathbb{Z} \backslash\{(n-1), n,(n+1)\} \tag{3.27}
\end{equation*}
$$

Let $v \in X$ and define $M=D_{1} G(B, \delta)-D_{1} G(0,0), w=M v$. Then for every $n \in \mathbb{Z}$ we have

$$
\begin{align*}
w_{n} & =\sum_{m \in \mathbb{Z}} M_{n, m} v_{m}=\sum_{m \in \mathbb{Z}}\left(\frac{\partial G_{n}}{\partial B_{m}}(B, \delta)-\frac{\partial G_{n}}{\partial B_{m}}(0,0)\right) v_{m}  \tag{3.28}\\
& =\delta(\Delta v)_{n}+2 \gamma \tanh \frac{\kappa}{2} \sum_{m \in \mathbb{Z}} \rho^{|m-n|}\left(\left|B_{m}+\widetilde{A}_{m}\right|^{2}-\left|\widetilde{A}_{m}\right|^{2}\right) v_{n}  \tag{3.29}\\
& +4 \gamma \tanh \frac{\kappa}{2} \sum_{m \in \mathbb{Z}} \rho^{|m-n|}\left(\left|B_{m}+\widetilde{A}_{m}\right|\left(B_{n}+\widetilde{A}_{n}\right)-\left|\widetilde{A}_{m}\right| \widetilde{A}_{n}\right) v_{m}  \tag{3.30}\\
& =\delta(\Delta v)_{n}+2|\gamma| \tanh \frac{\kappa}{2}\left(K_{n}+2 J_{n}\right), \tag{3.31}
\end{align*}
$$

where

$$
\begin{gather*}
K_{n}=\sum_{m \in \mathbb{Z}} \rho^{|m-n|}\left(\left|B_{m}+\widetilde{A}_{m}\right|^{2}-\left|\widetilde{A}_{m}\right|^{2}\right) v_{n}, \quad \forall n \in \mathbb{Z},  \tag{3.32}\\
J_{n}=\sum_{m \in \mathbb{Z}} \rho^{|m-n|}\left(\left|B_{m}+\widetilde{A}_{m}\right|\left(B_{n}+\widetilde{A}_{n}\right)-\left|\widetilde{A}_{m}\right| \widetilde{A}_{n}\right) v_{m}, \quad \forall n \in \mathbb{Z} . \tag{3.33}
\end{gather*}
$$

Hence

$$
\begin{equation*}
\|w\|_{X} \leq|\delta|\|\Delta\|_{X \times X}\|v\|_{X}+2|\gamma| \tanh \frac{\kappa}{2}\|K\|_{X}+4|\gamma| \tanh \frac{\kappa}{2}\|J\|_{X} \tag{3.34}
\end{equation*}
$$

Observe that by (3.32)

$$
\left|K_{n}\right| \leq \sum_{m \in \mathbb{Z}} \rho^{|m-n|}\left(B_{m}^{2}+2\left|B_{m}\right|\left|\widetilde{A}_{m}\right|\right)\left|v_{n}\right|,
$$

so that

$$
\begin{align*}
\|K\|_{X}^{2} & \leq \sum_{n \in \mathbb{Z}}\left|\sum_{m \in \mathbb{Z}} \rho^{|m-n|}\left(B_{m}^{2}+2\left|B_{m} \| \widetilde{A}_{m}\right|\right)\right| v_{n}| |^{2}  \tag{3.35}\\
& \leq \sum_{n \in \mathbb{Z}}\|B\|_{\infty}\left(\|B\|_{\infty}+\|\widetilde{A}\|_{\infty}\right)\left|\sum_{m \in \mathbb{Z}} \rho^{|m-n|}\right| v_{n}| |^{2}  \tag{3.36}\\
& \leq\|B\|_{\infty}\left(\|B\|_{\infty}+\|\widetilde{A}\|_{\infty}\right)\left\|\rho^{|\cdot|}\right\|_{1}^{2}\|v\|_{X}^{2} \tag{3.37}
\end{align*}
$$

using Young's inequality (3.12) with $u=\rho^{|\cdot|}, r=2, p=1$ and $q=2$. On the other hand, by (3.33) observe that

$$
\begin{equation*}
\left|J_{n}\right| \leq \sum_{m \in \mathbb{Z}} \rho^{|m-n|}\left(\left|B_{m}\right|\left|B_{n}\right|+\left|B_{m}\right|\left|\widetilde{A}_{n}\right|+\left|B_{n}\right|\left|\widetilde{A}_{m}\right|\right)\left|v_{m}\right|, \tag{3.38}
\end{equation*}
$$

so that

$$
\begin{align*}
\|J\|_{X}^{2} & \leq \sum_{n \in \mathbb{Z}}\left|\sum _ { m \in \mathbb { Z } } \rho ^ { | m - n | } \left(\left|B _ { m } \left\|B _ { n } \left|+\left|B _ { m } \left\|\left.\widetilde{A}_{n}\left|+\left|B_{n} \| \widetilde{A}_{m}\right|\right)\left|v_{m}\right|\right|^{2}\right.\right.\right.\right.\right.\right.\right.  \tag{3.39}\\
& \leq \sum_{n \in \mathbb{Z}}\|B\|_{\infty}\left(\|B\|_{\infty}+2\|\widetilde{A}\|_{\infty}\right)\left|\sum_{m \in \mathbb{Z}} \rho^{|m-n|}\right| v_{m}| |^{2}  \tag{3.40}\\
& \leq\|B\|_{\infty}\left(\|B\|_{\infty}+2\|\widetilde{A}\|_{\infty}\right)\left\|\rho^{|\cdot|}\right\|_{1}^{2}\|v\|_{X}^{2} \tag{3.41}
\end{align*}
$$

using again Young's inequality.
Using (3.35) and (3.40) in (3.34) we have that

$$
\begin{equation*}
\|w\| \leq\left(|\delta|\|\Delta\|_{X \times X}+6|\gamma| \tanh \frac{\kappa}{2}\|B\|\left(\|B\|+2\|\widetilde{A}\|_{\infty}\right)\left\|\rho^{|\cdot|}\right\|_{1}\right)\|v\| . \tag{3.42}
\end{equation*}
$$

By (3.29), (3.42) we thus have (3.24) with $\beta$ that vanishes as $(B, \delta) \rightarrow 0 \in X \times \mathbb{R}$, as required.

## 4. Spectral analysis of 1-peak and shelf-type breathers

In this section we present results on the spectra of the linearization around some of the breather solutions considered in the previous sections and in [BCMP15]. We will consider the one-peak breather solutions, related to the minimizers of section 2, and shelf-type breather of the local and nonlocal equation, studied in section 3 .

The linear stability analysis uses the fact that breather solutions are relative equilibria of (2.1) with respect to the action of global phase change. Using the variable $v$ defined by $u=e^{-i \omega t} v$, Hamilton's equation (2.2) becomes

$$
\begin{equation*}
\dot{v}_{n}=-i \frac{\partial H_{\omega}}{\partial v_{n}^{*}}, \quad n \in Z, \quad \text { with } \quad H_{\omega}=H-\omega P . \tag{4.1}
\end{equation*}
$$

A breather solution $u=e^{-i \omega t} A$ of (2.1) is a fixed point $v=A$ of (4.1).
Let $z=[q, p]^{T} \in X \times X$, with $z_{n}=\left[q_{n}, p_{n}\right]^{T}, q_{n}=\operatorname{Re} v_{n}, p_{n}=\operatorname{Im} v_{n}, n \in \mathbb{Z}$. Then (4.1) is also written as

$$
\begin{equation*}
\dot{z}=J \nabla h_{\omega}, \quad \text { with } \quad h_{\omega}=\frac{1}{2} H_{\omega}, \tag{4.2}
\end{equation*}
$$

and $(J z)_{n}=-\left[p_{n}, q_{n}\right]^{T}$, i.e. $J$ is the standard symplectic operator in $X \times X$.
The linearization at a fixed point $A$ of (4.1) is

$$
\begin{equation*}
\dot{z}=J \mathcal{H} z, \quad \text { with } \quad \mathcal{H}=\nabla^{2} h_{\omega}(A), \tag{4.3}
\end{equation*}
$$

i.e. $\mathcal{H}$ is the Hessian of $h_{\omega}$ at $A$.

To simplify the notation we set $\gamma=-1$ in this section. This does not affect the generality of the results.

We calculate $J \mathcal{H}$ explicitly as

$$
J=\left[\begin{array}{cc}
0 & I  \tag{4.4}\\
-I & 0
\end{array}\right], \quad \mathcal{H}=\left[\begin{array}{cc}
L_{+} & 0 \\
0 & L_{-}
\end{array}\right],
$$

where

$$
\begin{equation*}
L_{-}=-\omega I-\delta \Delta+2 \mathcal{A}, \quad L_{+}=-\omega I-\delta \Delta+2 \mathcal{A}+4 \mathcal{M} \tag{4.5}
\end{equation*}
$$

and $\mathcal{A}, \mathcal{M}$ are linear operators in $X$ defined by

$$
\begin{gather*}
\mathcal{A}(n, k)=\tanh \frac{\kappa}{2}\left(\sum_{m \in \mathbf{Z}} e^{-\kappa|m-n|} A_{m}^{2}\right) \delta_{n, k}, \quad n, \kappa \in \mathbb{Z},  \tag{4.6}\\
\mathcal{M}(n, k)=\tanh \frac{\kappa}{2} e^{-\kappa|n-k|} A_{k} A_{n}, \quad n, k \in \mathbb{Z}, \tag{4.7}
\end{gather*}
$$

with $\delta_{n, k}$ the Kronecker delta.
The linearization around breathers of the local cubic DNLS is obtained by taking the limit $\kappa \rightarrow \infty$ in (4.6), (4.7).

The above calculations are formal. In the case $A \in l^{\infty}(\mathbb{Z}, \mathbb{R})$, the operators $\mathcal{A}, \mathcal{M}$, $L_{-}, L_{+}$are bounded operators in $X$, while $J, \mathcal{H}$, and $J \mathcal{H}$ are bounded operators in $X \times X$. We also see that $L_{-}, L_{+}$and $\mathcal{H}$ are symmetric. Similar statements apply to the local case.

Let $X=l_{2}(\mathbb{Z}, \mathbb{R})$ with the inner product $(u, v)=\sum_{n \in \mathbb{Z}} u_{n} v_{n}$. The norm is denoted by $\|\cdot\|$. Let $X_{c}$ be the complexification of $X$, with the inner product $(u, v)_{c}=\sum_{n \in \mathbb{Z}} u_{n} v_{n}^{*}$.

Let $Y=l_{2}(\mathbb{Z}, \mathbb{C})$ viewed as a real space with the (real) inner product $\langle u, v\rangle=$ $\sum_{\mathbb{Z}}\left[\left(\operatorname{Re} u_{n}\right)\left(\operatorname{Re} v_{n}\right)+\left(\operatorname{Im} u_{n}\right)\left(\operatorname{Im} v_{n}\right)\right]$. Let $Y_{c}=l_{2}\left(\mathbb{Z}, \mathbb{C}^{2}\right)$ with the inner product $\langle u, v\rangle_{c}=$ $\sum_{\mathbb{Z}}\left[\left(\operatorname{Re} u_{n}\right)\left(\operatorname{Re} v_{n}\right)^{*}+\left(\operatorname{Im} u_{n}\right)\left(\operatorname{Im} v_{n}\right)^{*}\right] . \quad Y_{c}$ is the complexification of $Y$. Note that $Y=X \times X$, and $Y_{c}=X_{c} \times X_{c}$. The notation for $X, Y$ is that of the previous sections, where we only used the Banach space structure.

Let $M$ be a bounded operator in $Y_{c}$, or $X_{c}$. We denote the spectrum of $M$ by $\sigma(M)$. The essential spectrum $\sigma_{e}(M)$ of $M$ is defined as in [KP13], p.29, namely $\lambda \in \sigma(M)$ belongs to $\sigma_{e}(M)$ if $M-\lambda I$ either fails to be Fredholm, or is Fredholm with nonzero index. The definition of [Kato76], p.243, requires a weaker condition, namely $M-\lambda I$ not semi-Fredholm.

Let $\mathbf{X}$ be a real or complex Hilbert space, with norm $\|\cdot\| \mathbf{X}$, and let $\left\{\varepsilon_{j}\right\}_{j=1}^{\infty}$ be an orthonormal basis of $\mathbf{X}$. A bounded operator $T: \mathbf{X} \rightarrow \mathbf{X}$ is Hilbert-Schmidt in $\mathbf{X}$ if $\sum_{j=1}^{\infty}\left\|T \varepsilon_{j}\right\|_{\mathbf{X}}^{2}<\infty$ (this is independent of the choice of the basis). Recall that a Hilbert-Schmidt operator is compact.

### 4.1. Stability of one-peak breathers for small linear coupling

We will consider the one-peak breather solutions, defined below. These are solutions of (2.4) with $A \in X$, and start by showing some facts on arbitrary such solutions.
Lemma 4.1. Let $A \in X$. Then the operators $\mathcal{A}, \mathcal{M}$ defined in (4.6), (4.7) respectively are Hilbert-Schmidt in $X, X_{c}$.

Proof. We use the basis $\left\{\hat{e}_{n}\right\}_{n \in \mathbb{Z}}$ of $X$ and $X_{c}$. Both $\mathcal{A}, \mathcal{M}$ map $X \subset X_{c}$ to $X$, and it is enough to show the statement for $X$. First, by (4.6)

$$
\begin{equation*}
\sum_{n \in \mathbb{Z}}\left\|\mathcal{A} \hat{e}_{n}\right\|^{2}=\tanh ^{2} \frac{\kappa}{2} \sum_{n \in \mathbb{Z}}\left|\sum_{m \in \mathbb{Z}} e^{-\kappa|m-n|} A_{m}^{2}\right|^{2}=\tanh ^{2} \frac{\kappa}{2}\left\|g * A^{2}\right\|_{l^{2}}^{2}, \tag{4.8}
\end{equation*}
$$

where $g_{n}=e^{-\kappa|n|}$, and $\left(A^{2}\right)_{n}=A_{n}^{2}, \forall n \in \mathbb{Z}$. By the discrete Young's inequality (3.12) we have

$$
\begin{equation*}
\left\|g * A^{2}\right\|_{l^{2}} \leq\|g\|_{l^{1}}\left\|A^{2}\right\|_{l^{2}} \tag{4.9}
\end{equation*}
$$

Since $\|A\|_{l \infty} \leq\|A\|$, both quantities are finite, and (4.8), (4.9) imply that $\mathcal{A}$ is HilbertSchmidt in $X$. Similarly, by (4.7) we have

$$
\begin{equation*}
\sum_{\lambda \in \mathbb{Z}}\left\|\mathcal{M} \hat{e}_{\lambda}\right\|^{2}=\tanh ^{2} \frac{\kappa}{2} \sum_{\lambda \in \mathbb{Z}} A_{\lambda}^{2} \sum_{n \in \mathbb{Z}}\left|e^{-\kappa|\lambda-n|} A_{n}\right|^{2} \leq\|A\|_{l^{2}}^{4}, \tag{4.10}
\end{equation*}
$$

therefore $\mathcal{M}$ is also Hilbert-Schmidt in $X$. $\square$
Lemma 4.1 is applied to fixed points $A \in X$ of (4.1). Define $\mathcal{H}_{0}$ as

$$
\mathcal{H}_{0}=\left[\begin{array}{cc}
-\omega I-\delta \Delta & 0  \tag{4.11}\\
0 & -\omega I-\delta \Delta
\end{array}\right]
$$

Clearly $\mathcal{H}_{0}$ is a bounded operator in $Y_{c}$, moreover $J \mathcal{H}$ is a compact perturbation of $J \mathcal{H}_{0}$ :
Proposition 4.2. Let $A \in X$, and let $J \mathcal{H}: Y \rightarrow Y$ be as in (4.4)-(4.7). Then (i) $\sigma_{e}(J \mathcal{H})=\sigma_{e}\left(J \mathcal{H}_{0}\right)$, and (ii) $\sigma_{e}\left(J \mathcal{H}_{0}\right)$ consists of $z \in \mathbf{C}$ with $\operatorname{Re} z=0$, $\operatorname{Imz} \in[-\omega,-\omega+$ $4 \delta] \cup[\omega-4 \delta, \omega]$, for $\delta>0$, and $\operatorname{Rez}=0$, $\operatorname{Imz} \in[-\omega+4 \delta,-\omega] \cup[\omega, \omega-4 \delta]$, if $\delta<0$.

Proof of Part (i). To show that $\sigma_{e}(J \mathcal{H})=\sigma_{e}\left(J \mathcal{H}_{0}\right)$ it is enough to show that $J \mathcal{H}-$ $J \mathcal{H}_{0}$ is Hilbert-Schmidt in $Y$. We have

$$
J \mathcal{H}-J \mathcal{H}_{0}=\left[\begin{array}{cc}
0 & 2 \mathcal{A}  \tag{4.12}\\
-2 \mathcal{A}+4 \mathcal{M} & 0
\end{array}\right] .
$$

The Hilbert-Schmidt property in $Y=X \times X$ follows easily from Lemma 4.1. We can use the basis vectors $\hat{e}_{n}^{j}, n \in \mathbb{Z}, j=1,2$, defined as $\hat{e}_{n}^{1}(k)=\delta_{n, k}[1,0]^{T}, \hat{e}_{n}^{2}(k)=\delta_{n, k}[0,1]^{T}$, for all $n, k \in \mathbb{Z}$. Part (ii) is shown in the next subsection.

We now consider the stability of the 1 -peak solution studied numerically in [BCMP15]. Consider a solution $A \in X$ of (2.4) with $\delta=0$ that satisfies $A_{n}=0$, for all $n \in \mathbb{Z} \backslash\left\{n_{0}\right\}$, and $A_{n_{0}}>0$. The parameter $\kappa>0$ is fixed and arbitrary. Such a solution is unique, and by Theorem 2.4 of [BCMP15] it can be continued uniquely to a $C^{0}$ family of solutions $A(\delta), \delta \in\left(-\delta_{0}, \delta_{0}\right)$, of (2.4). We refer to these solutions as 1-peak breathers. Remark 4.3. Note that the continuation result in [BCMP15] is for solutions with $\|A(\delta)\|_{X}$ fixed. Then $\omega$ is also a function of $\delta$, and is not known a-priori. The continuation result for fixed $\omega$ has a very similar proof, this is seen in the previous section.

For the $\delta=0$ breather $A(0), J \mathcal{H}$ is block diagonal with $2 \times 2$ blocks, each block corresponding to an integer $k \in \mathbb{Z}$. The block corresponding to $k=n_{0}$ has a double eigenvalue $\lambda_{n_{0}}=0$. The other blocks have eigenvalues

$$
\begin{equation*}
\pm \lambda_{k}= \pm i \omega\left(1-e^{\kappa\left|k-n_{0}\right|}\right), \quad \forall k \in \mathbb{Z} \backslash\left\{n_{0}\right\} \tag{4.13}
\end{equation*}
$$

with corresponding eigenfunctions $w_{k}^{ \pm} \in Y_{c}$, given by

$$
\begin{equation*}
w_{k}^{ \pm}(n)=\delta_{k, n}[1, \mp i]^{T}, \quad n \in \mathbb{Z} . \tag{4.14}
\end{equation*}
$$

Thus $\sigma(J \mathcal{H})$ consists of a double zero eigenvalue, an infinite number of pairs of imaginary eigenvalues $\pm \lambda_{k}, \kappa>n_{0}$, of multiplicity 2 , and their accumulation points $\pm i \omega$ (belonging to $\sigma_{e}(J \mathcal{H})$ ).

Considering now the linearized operator $J \mathcal{H}$ corresponding to a 1-peak breather $A(\delta)$, we have by Proposition 4.2 that $\sigma_{e}(J \mathcal{H})=\sigma\left(J \mathcal{H}_{0}\right)$. The persistence of the double zero eigenvalue is a well known fact, it follows from the fact that $L_{-} A(\delta)=0$, the continuity of finite rank spectral projections, and the symmetry of spectra of symplectic operators. We now consider the persistence of the internal mode eigenvalues $\lambda_{k}$ of (4.13). The linearized operator (4.4)-(4.7) corresponding to the 1-peak breather $A(\delta)$, $\delta \in\left(-\delta_{0}, \delta_{0}\right)$, will be denoted by $J \mathcal{H}(\delta)$.
Proposition 4.4. Let $k>n_{0}$, and consider the pair of imaginary eigenvalues $\pm \lambda=$ $\pm \lambda_{k}$ of the $J \mathcal{H}(0)$, as in (4.13). Then there exists a $\delta_{0, k}>0$ for which the operators $J \mathcal{H}(\delta): Y \rightarrow Y, \delta \in\left(-\delta_{0, k}, \delta_{0, k}\right)$, have two eigenvalues $\pm \lambda_{j}(\delta) \in i \mathbb{R}, j=1,2$. The families $\pm \lambda_{j}(\delta), j=1,2$ are continuous in $\delta$ and satisfy $\pm \lambda_{j}(0)= \pm \lambda$.

Proof. Fix $k>n_{0}$. By the continuity of finite dimensional spectral projections and the corresponding eigenvalues under perturbations by bounded operators, see [Kato76], the families $\pm \lambda_{j}(\delta), j=1,2$ are continuous in $\delta$, with $\pm \lambda(0)= \pm \lambda= \pm \lambda_{k}$, as stated, assuming that $\delta$ belongs to some interval $\left(-\delta_{0, k}, \delta_{0, k}\right)$. Consider the eigenvalue $\lambda(0)$ of $J \mathcal{H}(0)$. By the symmetry of spectra of symplectic operators the two continued eigenvalues of $J \mathcal{H}(\delta)$ are either on the imaginary axis, or are off the imaginary axis and satisfy $\lambda_{2}(\delta)=-\lambda_{1}^{*}(\delta)$. To eliminate the second scenario we use the Krein dichotomy argument, see [KP13], p. 181, (7.14). In particular, if the $\lambda_{j}(\delta)$ have nonzero real parts then the corresponding eigenvectors $v_{j}(\delta)$ must satisfy

$$
\begin{equation*}
\left\langle\mathcal{H}(\delta) v_{j}(\delta), v_{j}(\delta)\right\rangle_{c}=0, \quad \forall \delta \in\left(-\delta_{0, k}, \delta_{0, k}\right), \quad j=1,2 . \tag{4.15}
\end{equation*}
$$

On the other hand, using explicit expressions for the eigenvectors $v_{j}(0), j=1,2$, i.e. $v_{1}(0)=w_{k}^{+}, v_{2}(0)=w_{m}^{+}$, with $m<n_{0}$ satisfying $n_{0}-m=k-n_{0}$, and $w_{k}^{+}, w_{m}^{+}$as in (4.14), we compute that

$$
\begin{equation*}
\left\langle\mathcal{H}(0) v_{j}(0), v_{j}(0)\right\rangle_{c}=-2 \omega\left(1-e^{\kappa\left|k-n_{0}\right|}\right) \neq 0, \quad j=1,2, \tag{4.16}
\end{equation*}
$$

contradicting (4.15) at $\delta=0$.
We expect that for $\delta, \gamma>0, \delta$ sufficiently small, the 1 -peak breathers obtained by continuation should coincide with the global minimizers of Section 2. Nevertheless at the moment we can prove the following weaker statement.
Proposition 4.5. Let $\lambda>0$, and fix $n_{0} \in \mathbb{Z}$. Consider the continuous one-parameter curve of $A(\delta) \in X, \delta \in\left[0, \delta_{0}\right]$, of solutions of (2.4) with $A(0)$ the 1-peak breather solution supported at $n=n_{0}$, and $\|A(\delta)\|=\lambda, \forall \delta \in\left[0, \delta_{0}\right]$. Then there is a sequence of $\left\{\delta_{m}\right\}_{m \in \mathbb{Z}^{+}} \subset\left[0, \delta_{0}\right]$, that converges to $\delta=0$, and such that $A\left(\delta_{m}\right)$ minimizes $H$ of (2.3) over all configurations $u \in X$ with $P(u)=\lambda$.

Proof. Fix $\lambda>0$. Let $H_{\delta}=H=\delta H_{2}-H_{4}$, where $H_{2}, V$ are the quadratic, and quartic parts of $H$ respectively, see (2.3). For each $\delta \in\left[0, \delta_{0}\right]$, let $u_{*}^{\delta} \in X$ be a minimizer of $H_{\delta}$ at power $P=\lambda$. For $\delta^{\prime}, \delta \in\left[0, \delta_{0}\right]$, we have that $\delta<\delta^{\prime}$ implies $H_{\delta}\left(u_{*}^{\delta}\right) \leq H_{\delta}\left(u_{*}^{\delta^{\prime}}\right) \leq$ $H_{\delta^{\prime}}\left(u_{*}^{\delta^{\prime}}\right)$. Also, we easily check that $\delta^{\prime} \searrow \delta$ implies $H_{\delta^{\prime}}\left(u_{*}^{\delta^{\prime}}\right) \rightarrow H_{\delta}\left(u_{*}^{\delta}\right)$. It follows that any sequence of points $u_{*}^{\delta}, \delta \in\left[0, \delta_{0}\right]$, that converges to $u_{*}^{0}$ is a minimizing sequence for $H_{0}=V$, and therefore has a convergent subsequence, up to translations, by the arguments of Theorem 2.1. On the other hand, the critical points of $V$ are solutions of (3.6). The support of minimizers with finite support consists of sets of consecutive sites by Theorem 2.2. Such solutions were described in [BCMP15], where we saw that all solutions with support consisting of more that three sites must have their maxima at interfaces. Then the only solution having the properties of Theorem 2.2 is the one supported on one site. We also observe that by Young's inequality, the right hand side of (3.6) maps $X$ to $X$. Thus there are no solutions with infinite support. The points of the minimizing subsequence are solutions of (2.4), and must therefore come arbitrarily close to the breather $A(0)$, after suitable translations. By the uniqueness of the continuation, Theorem 2.4, [BCMP15], these points must eventually belong to the branch $A(\delta)$ of solutions obtained for some interval $\left[0, \delta_{0}\right]$.

### 4.2. Essential spectrum of linearization around shelf-type breather: local case

We now consider the essential spectrum of the linearization around the shelf solutions of section 3 . We first consider the shelf solutions of the local (power) nonlinearity system. The main result is Theorem 4.11 below. It follows from an analysis of the piecewise linear discrete dynamical system defined by the linear operator. We give a fairly complete presentation of the theory. In the nonlocal case we have a weaker result, Theorem 4.13 below, using similar computations. The theory however is not complete.

In the case of the power nonlinearity system, we consider solutions obtained by continuation of the $\delta=0$ breather

$$
\widetilde{A}_{n}=\left\{\begin{array}{c}
\alpha, \text { if } n \leq 0  \tag{4.17}\\
0, \text { if } n \geq 1,
\end{array}\right.
$$

with $\alpha=\sqrt{\frac{\omega}{2}}$. This is a special case of the configuration of (3.3). The spectral results below are the same for the general case.

For $\omega$ fixed, let $B_{n}, n \in \mathbb{Z}$, be a solution of Theorem 3.2 and let

$$
\begin{equation*}
A_{n}=\widetilde{A}_{n}+B_{n}, \quad n \in \mathbb{Z} \tag{4.18}
\end{equation*}
$$

Let $\mathcal{L}_{\infty, \delta}=J \mathcal{H}$, with $J \mathcal{H}$ as in (4.4), $L_{-}, L_{+}$as in (4.5), $\mathcal{A}, \mathcal{M}$ as in (4.6), (4.7) with $A_{n}$ as in (4.18), $\widetilde{A}_{n}$ as in (4.17), $\forall n \in \mathbb{Z}$. Operator $\mathcal{L}_{\infty, \delta}$ is therefore the formal linearization around the shelf solution of Theorem 3.2.

Also let $\mathcal{L}=J \mathcal{H}$, with $J \mathcal{H}$ as in (4.4), with $L_{-}, L_{+}$as in (4.5), and

$$
\begin{equation*}
\mathcal{A}(n, k)=\mathcal{A}_{0}(n, k)=\widetilde{A}_{n}^{2} \delta_{n, k}, \quad \mathcal{M}(n, k)=\mathcal{M}_{0}(n, k)=\widetilde{A}_{n}^{2} \delta_{n, k}, \tag{4.19}
\end{equation*}
$$

where $\widetilde{A}_{n}$ is as in (4.17).
Proposition 4.6. The operator $\mathcal{L}_{\infty, \delta}$ is a compact perturbation of $\mathcal{L}$, and therefore $\sigma_{e}\left(\mathcal{L}_{\infty, \delta}\right)=\sigma_{e}(\mathcal{L})$ in $Y$.

The prove the statement it is enough to show that $\mathcal{A}-\mathcal{A}_{0}, \mathcal{M}-\mathcal{M}_{0}$ are HilbertSchmidt. The proof of this fact is a simplified version of the proof of Lemma 4.12 for the nonlocal problem below, and is omitted.

To calculate the essential spectrum of $\mathcal{L}$ we study the equation $(\mathcal{L}-\lambda) z=0, z=$ $[q, p]^{T}$, with $z_{n}=\left[q_{n}, p_{n}\right]^{T}$ as a discrete dynamical system. The (spectral) parameter $\lambda$ varies over $\mathbb{C}$. By the definition of $\mathcal{L}$, with $\mathcal{A}, \mathcal{M}$ of (4.19), and $\widetilde{A}$ of (4.17) we have that $\mathcal{L} z=\lambda z$ is equivalent to

$$
\begin{align*}
& -\delta\left(p_{n+1}+p_{n-1}-2 p_{n}\right)=\lambda q_{n}, \quad-2 \omega q_{n}+\delta\left(q_{n+1}+q_{n-1}-2 q_{n}\right)=\lambda p_{n}, \quad n \leq 0  \tag{4.20}\\
& -\omega p_{n}-\delta\left(p_{n+1}+p_{n-1}-2 p_{n}\right)=\lambda q_{n}, \quad \omega q_{n}+\delta\left(q_{n+1}+q_{n-1}-2 q_{n}\right)=\lambda p_{n}, \quad n \geq 1 \tag{4.21}
\end{align*}
$$

Finite difference equation (4.20), (4.21) is furthermore equivalent to a first order system $y_{n+1}=\mathcal{F}\left(y_{n}\right)$ for a sequence of vectors $y_{n}=\left[y_{n}(1), \ldots, y_{n}(4)\right]^{T} \in \mathbb{C}^{4}, n \in \mathbb{Z}$. This formulation uses the identification

$$
\begin{equation*}
\left[y_{n}(1), y_{n}(2), y_{n}(3), y_{n}(4)\right]^{T}=\left[q_{n-1}, q_{n}, p_{n-1}, p_{n}\right]^{T} \tag{4.22}
\end{equation*}
$$

By (4.22), $n \leq 0$, (4.20) is then equivalent to

$$
y_{n+1}=M_{-}(\lambda) y_{n}, \quad n \leq 0, \quad \text { with } \quad M_{-}(\lambda)=\left[\begin{array}{cccc}
0 & 1 & 0 & 0  \tag{4.23}\\
-1 & 2 \frac{\omega}{\delta}+2 & 0 & \frac{\lambda}{\delta} \\
0 & 0 & 0 & 1 \\
0 & -\frac{\lambda}{\delta} & -1 & 2
\end{array}\right] .
$$

Similarly by (4.22), $n \geq 1$, (4.21) is equivalent to

$$
y_{n+1}=M_{+}(\lambda) y_{n}, \quad n \geq 1, \quad \text { with } \quad M_{+}(\lambda)=\left[\begin{array}{cccc}
0 & 1 & 0 & 0  \tag{4.24}\\
-1-\frac{\omega}{\delta}+2 & 0 & \frac{\lambda}{\delta} \\
0 & 0 & 0 & 1 \\
0 & -\frac{\lambda}{\delta} & -1 & -\frac{\omega}{\delta}+2
\end{array}\right] .
$$

Equations (4.23), (4.24) define a piecewise linear discrete-time dynamical system that may be written as $y_{n+1}=M(\lambda) y_{n}$, where $M(\lambda)=M_{-}(\lambda)$ if $n \leq 0$, and $M(\lambda)=$ $M_{+}(\lambda)$ if $n \geq 1$.
Proposition 4.7. Given any $m \in \mathbb{Z}, a \in \mathbb{C}^{4}$, there exists a unique $y_{n}, n \in \mathbb{Z}$, satisfying (4.23), (4.24) and $y_{m}=a$.

Proof. We check first that the $M_{ \pm}(\lambda)$ are invertible, $\forall \lambda \in \mathbb{C}$. By (4.23), (4.24)

$$
\begin{gather*}
P_{-}(r ; \lambda)=\operatorname{det}\left(M_{-}(\lambda)-r I\right)=\left[(r-1)^{2}-2 r \frac{\omega}{\delta}\right](r-1)^{2}+r^{2} \frac{\lambda^{2}}{\delta^{2}}  \tag{4.25}\\
P_{+}(r ; \lambda)=\operatorname{det}\left(M_{+}(\lambda)-r I\right)=\left[-r\left(2-\frac{\omega}{\delta}-r\right)+1\right]^{2}+r^{2} \frac{\lambda^{2}}{\delta^{2}} \tag{4.26}
\end{gather*}
$$

We see that $P_{ \pm}(r, \lambda)=0, r=0$ leads to a contradiction, $\forall \lambda \in \mathbb{C}$. Thus the $M_{ \pm}(\lambda)$ are invertible, $\forall \lambda \in \mathbb{C}$. This fact allows us to iterate (4.23), (4.24) backwards as well, writing

$$
\begin{align*}
& {\left[M_{-}(\lambda)\right]^{-1} y_{n+1}=y_{n}, \quad n \leq 0,}  \tag{4.27}\\
& {\left[M_{+}(\lambda)\right]^{-1} y_{n+1}=y_{n}, \quad n \geq 1 .} \tag{4.28}
\end{align*}
$$

We can thus iterate both forwards and backwards in each segment. We check that (4.23), (4.24), (4.27), and (4.28) allow us to cross the interface separating the $M_{ \pm}(\lambda)$ segments in a unique way, in either direction.

We examine the dynamics of each of the systems (4.23), (4.24), particularly its dependence on the parameter $\lambda$. In both cases we can start with some $y_{1}$ and consider the iteration backwards, and forwards respectively. Since each system is linear this amounts to a spectral analysis of the matrices $M_{ \pm}(\lambda)$. The results of this analysis are then interpreted in operator-analytic terms for the resolvent equation $(\mathcal{L}-\lambda) z=f$, $f=[g, h]^{T} \in X \times X \subset Y_{c}$.

Consider a linear system $x_{n+1}=T x_{n}$ in $\mathbb{C}^{N}$, where $T$ an invertible $N \times N$ matrix, and $n$ belongs to a set of indices $\mathcal{J}$ that is either $\mathbb{Z}$, or the set of all integers above or below some $m \in \mathbb{Z}$. The system is hyperbolic if there exist subspaces $E^{s}, E^{u}$ of $\mathbb{C}^{N}$, $\operatorname{dim}\left(E^{s}\right)+\operatorname{dim}\left(E^{u}\right)=N$ that are invariant under $T$, and for which there exist $C>0$, $\varrho>1$ such that

$$
\begin{equation*}
\left\|T^{n} P^{s} v\right\|_{\mathbb{C}^{N}} \leq C \varrho^{-n}\left\|P^{s} v\right\|_{\mathbb{C}^{N}}, \quad\left\|T^{-n} P^{u} v\right\|_{\mathbb{C}^{N}} \leq C \varrho^{-n}\left\|P^{u} v\right\|_{\mathbb{C}^{N}}, \quad \forall n \in \mathcal{J} \tag{4.29}
\end{equation*}
$$

where $P^{s}, P^{u}$ are respective projections to $E^{s}, E^{u}$ (and commute with $T$ ). The subspaces $E^{s}, E^{u}$ are referred to as stable and unstable subspaces respectively. Hyperbolicity can be often verified by calculating the eigenvalues of $T$. The dimensions of $E^{s}$, $E^{u}$ are determined be the multiplicities of eigenvalues inside and outside the unit circle respectively.

In the present problem hyperbolicity will be examined for each of the systems (4.23), (4.24). The corresponding stable and unstable subspaces will be denoted by $E_{ \pm}^{s}, E_{ \pm}^{u}$. Note that the $E_{-}^{s}, E_{-}^{u}$ for (4.23) refer to the forward iteration of $M_{-}(\lambda)$.

Lemmas 4.8, 4.9 connect the analysis of the eigenvalues of the $M_{ \pm}(\lambda)$ to the of spectrum of $\mathcal{L}$. The proofs are at the end of the subsection.
Lemma 4.8. Let $\lambda \in \mathbb{C}$ be such that the systems (4.23), (4.24) defined by $M_{-}(\lambda), M_{+}(\lambda)$ respectively are both hyperbolic, with $\operatorname{dim}\left(E_{-}^{u}\right)+\operatorname{dim}\left(E_{+}^{s}\right)=N, E_{-}^{u} \cap E_{+}^{s}=\{0\}$. Then $\lambda$ belongs to the residual set of $\mathcal{L}$ in $Y$.
Lemma 4.9. Let $\lambda \in \mathbb{C}$, and suppose that one of $M_{-}(\lambda), M_{+}(\lambda)$ has a semisimple eigenvalue on the unit circle. Then $\lambda$ belongs to the essential spectrum of $\mathcal{L}$ in $Y$.

Define the sets

$$
\begin{gather*}
\left.B_{1}=\left\{-\left(\omega-4 \delta \sin ^{2} \frac{k}{2}\right)^{2}: \quad k \in \mathbb{R}\right)\right\},  \tag{4.30}\\
\left.B_{2}=\left\{-4 \delta \sin ^{2} \frac{k}{2}\left(2 \omega+4 \delta \sin ^{2} \frac{k}{2}\right): \quad k \in \mathbb{R}\right)\right\} . \tag{4.31}
\end{gather*}
$$

We will further assume that $B_{1}, B_{2}$ are disjoint. By (4.30), (4.31) this property is clearly satisfied for $|\delta|$ sufficiently small.
Lemma 4.10. Let $\lambda \in \mathbb{C}$. If $\lambda^{2} \notin B_{1} \cup B_{2}$ then $\lambda$ either belongs to the residual set of $\mathcal{L}: Y \rightarrow Y$, or is an eigenvalue of $\mathcal{L}$ of finite multiplicity.
Theorem 4.11. The essential spectrum of the operators $\mathcal{L}$ and $\mathcal{L}_{\infty, \kappa}$ in $Y$ consists of the set of $\lambda \in \mathbb{C}$ satisfying $\lambda^{2} \in B_{1} \cup B_{2}$.

The case $\lambda^{2} \in B_{1}$ yields two intervals of $\lambda$ on the imaginary axis, and coincides with the essential spectrum of the linearization around breathers that decay at infinity, as in Proposition 4.2 (its proof is completed below).

The case $\lambda^{2} \in B_{2}$ yields two intervals that are both on either the imaginary or the real axis, depending on $\delta$ and $\omega$. For instance, for $|\delta|<\frac{|\omega|}{2}$ we have stability for $\delta \omega>0$, and instability for $\delta \omega<0$. In both cases zero is one the endpoints of the intervals, and is included in $\sigma_{e}(\mathcal{L})$.

Proof of Lemma 4.10. The eigenvalues of $M_{ \pm}(\lambda)$ are the roots of the polynomials $\left.\operatorname{det}\left(M_{ \pm}(\lambda)-r I\right)\right)=P_{ \pm}(r ; \lambda)$, shown in (4.25), (4.26).

We can rewrite $P_{-}(r ; \lambda)=0$ as

$$
\begin{equation*}
\lambda^{2}=\delta\left(r+r^{-1}-2\right)\left[2 \omega-\delta\left(r+r^{-1}-2\right)\right] . \tag{4.32}
\end{equation*}
$$

If a root $r$ is on the unit circle, i.e. $r=e^{i k}$ for some real $k$, then $\lambda^{2}$ must belong to $B_{1}$. Moreover the roots of $P_{-}$come in pairs $r, r^{-1}$. If $\lambda^{2} \notin B_{1}$ we therefore have two roots inside and two roots outside the unit circle. The system defined by $M_{-}(\lambda)$ is hyperbolic, and $\operatorname{dim}\left(E_{-}^{s}\right)=\operatorname{dim}\left(E_{-}^{u}\right)=2$.

Analogously, $P_{+}(r ; \lambda)=0$ can be written as

$$
\begin{equation*}
\lambda^{2}=-\left[\delta^{2}\left(r+r^{-1}\right)^{2}+(\omega-2 \delta)^{2}+2 \delta\left(r+r^{-1}\right)\right] . \tag{4.33}
\end{equation*}
$$

It follows that if $r=e^{i k}$ for some real $k$, then $\lambda^{2}$ must belong to $B_{2}$. Similarly, since the roots of $P_{+}$come in pairs $r, r^{-1}$, if $\lambda^{2} \notin B_{2}$ then the system defined by $M_{+}(\lambda)$ is hyperbolic, and $\operatorname{dim}\left(E_{+}^{s}\right)=\operatorname{dim}\left(E_{+}^{u}\right)=2$.

Therefore $\lambda^{2} \notin B_{1} \cup B_{2}$ implies $\operatorname{dim}\left(E_{-}^{u}\right)+\operatorname{dim}\left(E_{+}^{s}\right)=4$. If in addition $E_{-}^{u} \cap E_{+}^{s}=$ $\{0\}$, Lemma 4.8 implies that $\lambda$ is in the residual set of $\mathcal{L}$. If otherwise $E_{-}^{u}=E_{+}^{s}$ then $(\mathcal{L}-\lambda) z=0$ has exponentially decaying solutions and is therefore an eigenvalue of $\mathcal{L}$. The set of these solutions is clearly two-dimensional.

Proof of Theorem 4.11. By the calculation in the proof of Lemma 4.10, if $\lambda^{2}$ belongs to $B_{1} \cup B_{2}$ then either $M_{-}(\lambda)$ or $M_{+}(\lambda)$ has an eigenvalue of the form $r=e^{i k}$ for some real $k$. By Lemma 4.9 then $\lambda$ belongs to the essential spectrum of $\mathcal{L}$. If, on the other hand, $\lambda^{2} \notin B_{1} \cup B_{2}$, Lemma 4.10 implies that $\lambda$ can not belong to $\sigma_{e}(\mathcal{L})$.

Proof of Proposition 4.2, Part (ii). Solving $J \mathcal{H}_{0} u=\lambda u$ explicitly i.e. using $u_{n}=e^{i k n}$, $\sigma_{e}\left(J \mathcal{H}_{0}\right)$ is as in the statement by Lemma 4.9. The result coincides with the part of the essential spectrum of $\mathcal{L}$ that comes from the set $B_{1}$ of (4.30).

We now prove Lemmas 4.8, and 4.9.
Proof of Lemma 4.8. The resolvent equation $(\mathcal{L}-\lambda) z=f, z=[p, q]^{T}, f=[g, h]^{T}$ in $X \times X \subset Y_{c}$ is equivalent to

$$
\begin{equation*}
L_{-} p=\lambda q+g, \quad-L_{+} q=\lambda p+h, \tag{4.34}
\end{equation*}
$$

and can be also written in the dynamical system form

$$
\begin{equation*}
y_{n+1}=M y_{n}+F_{n}, \quad \text { with } \quad F_{n}=\left[0, h_{n}, 0,-g_{n}\right]^{T}, \tag{4.35}
\end{equation*}
$$

where $M=M_{-}(\lambda)$ if $n \leq 0, M=M_{+}(\lambda)$ if $n \geq 1$. The solution of (4.35) is obtained by iterating forwards or backwards in time. We will examine forward iterates, from $n=1$ to $n>1$ arbitrarily large, similar considerations apply to backward iterates from $n=1$.

By (4.35) the solution for $n \geq 1$ is

$$
\begin{equation*}
y_{n+1}=M^{n} y_{1}+\sum_{j=1}^{n-1} M^{j} F_{n-j}, \quad \text { with } \quad M=M_{+}(\lambda) . \tag{4.36}
\end{equation*}
$$

Using the assumption of hyperbolicity, we can use projections $P^{s}, P^{u}$ to the stable and unstable subspaces respectively to write (4.36) as

$$
\begin{equation*}
y_{n+1}=M^{n} P^{s} y_{1}+\sum_{j=1}^{n-1} M^{j} P^{s} F_{n-j}+M^{n} P^{u} y_{1}+\sum_{j=1}^{n-1} M^{j} P^{u} F_{n-j}, \quad \forall n>1 \tag{4.37}
\end{equation*}
$$

We first claim that if $[g, h]^{T}$ belongs to $X_{c} \times X_{c}$ then the first two terms in the right hand side of (4.37) belong to $X_{c,+}^{4}=X_{c,+} \times X_{c,+} \times X_{c,+} \times X_{c,+}$, where $X_{c,+}$ is the space of configurations in $X_{c}$ that vanish for $n \leq 0$.

The claim is immediate for the first term, since $\left\|M^{n} P^{s} y_{1}\right\|_{\mathbb{C}^{4}}$ decays exponentialy in $n$ by the hyperbolicity assumption, see (4.29). Denote the second term by $G_{n}$, then by (4.37)

$$
\begin{equation*}
G_{n}=\sum_{j=1}^{n-1} M^{j} P^{s} F_{n-j}=\sum_{k=1}^{n} M^{n-k} P^{s} F_{k}, \tag{4.38}
\end{equation*}
$$

therefore by the hyperbolicity assumption, see (4.29),

$$
\begin{equation*}
\left\|G_{n}\right\|_{\mathbb{C}^{4}} \leq C \sum_{k=1}^{n} \rho_{2}^{-|n-k|}\left\|F_{k}\right\|_{\mathbb{C}^{4}} \leq C \sum_{k=1}^{\infty} \rho_{2}^{|n-k|}\left\|F_{k}\right\|_{\mathbb{C}^{4}} \tag{4.39}
\end{equation*}
$$

for some $C>0$, and $\rho_{2}>1$. Note that $[g, h]^{T} \in X_{c} \times X_{c}$ implies $F \in X_{c}^{4}$. Let $\tilde{f}_{k}=\left\|F_{k}\right\|_{\mathbb{C}^{4}}$ for $k \geq 1, \tilde{f}_{k}=0$ for $k \leq 0$, and define

$$
\begin{equation*}
\tilde{g}_{n}=C \sum_{k \in \mathbb{Z}} \rho_{2}^{-|n-k|} \tilde{f}_{k}, \quad n \in \mathbb{Z} \tag{4.40}
\end{equation*}
$$

Then $\tilde{g}_{n} \leq\left\|G_{n}\right\|_{\mathbb{C}^{4}}$ for all $n \geq 1$, and by (4.39), and Young's inequality (3.12),

$$
\begin{equation*}
\sum_{n=1}^{\infty}\left\|G_{n}\right\|_{\mathbb{C}^{4}}^{2} \leq\|\tilde{g}\|_{l^{2}}^{2} \leq C^{2}\left\|\rho_{2}^{-|\cdot|}\right\|_{l^{1}}^{2}\|\tilde{f}\|_{l^{2}}^{2}=C^{2}\left\|\rho_{2}^{-|\cdot|}\right\|_{l^{1}}^{2}\|F\|_{X_{c}^{4}}^{2} . \tag{4.41}
\end{equation*}
$$

Thus the second term of the right hand side of (4.37) belongs to $X_{c,+}^{4}$, as claimed.
Consider the third and fourth terms in the right hand side of (4.37). They can be written as $M^{n} B_{n}$ with

$$
\begin{equation*}
B_{n}=P^{u} y_{1}+\sum_{j=1}^{n-1} M^{-(n-j)} P^{u} F_{n-j} \tag{4.42}
\end{equation*}
$$

where $M^{-p}=\left(M^{-1}\right)^{p}$ if $p$ a positive integer. To have the sequence of $y_{n}$ of (4.37) in $X_{c,+}^{4}$ we need that $B_{n} \rightarrow 0$ as $n \rightarrow \infty$. By (4.42) we therefore need $P^{u} y_{1}=b_{+}$, with

$$
\begin{equation*}
b_{+}=-\lim _{n \rightarrow \infty} \sum_{j=1}^{n-1} M^{-(n-j)} P^{u} F_{n-j} \tag{4.43}
\end{equation*}
$$

We will check below that the limit exists, and that the solution of (4.36), $n \geq 1$, obtained for such $y_{1}$ belongs to $X_{c,+}^{4}$.

To show that the limit in (4.43) exists we must check that

$$
\begin{equation*}
I_{n}=\sum_{j=0}^{n-1} M^{-(n-j)} P^{u} F_{n-j}=\sum_{k=1}^{n} M^{-k} F_{k} \tag{4.44}
\end{equation*}
$$

converges. For $n>m$ we have

$$
\begin{equation*}
\left\|I_{n}-I_{m}\right\|_{\mathbb{C}^{4}} \leq \sum_{k=m+1}^{n}\left\|M^{-k} P^{u} F_{k}\right\|_{\mathbb{C}^{4}} \leq C \sum_{k=m+1}^{n} \rho_{2}^{k}\left\|F_{k}\right\|_{\mathbb{C}^{4}} \tag{4.45}
\end{equation*}
$$

for some $\rho_{2}<1, C>0$, by the hyperbolicity assumption. Then

$$
\begin{equation*}
\left\|I_{n}-I_{m}\right\|_{\mathbb{C}^{4}} \leq C\left(\sum_{k=m+1}^{n} \rho_{2}^{2 k}\right)^{\frac{1}{2}}\|F\|_{X^{4}} \leq C \rho_{2}^{m-1} C_{\rho_{2}}\|F\|_{X_{c}^{4}}, \tag{4.46}
\end{equation*}
$$

where $C_{\rho_{2}}$ is independent of $n, m$. Thus $I_{n}$ converges as $n \rightarrow \infty$, and $b_{+}$in (4.43) is well defined.

To check that the sum $J_{n}$ of the third and fourth terms in the right hand side of (4.37) obtained by choosing $y_{1}$ as above belongs to $X_{+}^{4}$, we note that for $n \geq 1$,

$$
\begin{align*}
J_{n} & =M^{n}\left[\lim _{k \rightarrow \infty} \sum_{j=0}^{k-1} M^{-(k-j)} P^{u} F_{k-j}-\sum_{j=0}^{n-1} M^{-(n-j)} P^{u} F_{k-j}\right]  \tag{4.47}\\
& =M^{n}\left[\sum_{l=1}^{+\infty} M^{-l} P^{u} F_{l}-\sum_{l=1}^{n} M^{-l} P^{u} F_{l}\right]  \tag{4.48}\\
& =\sum_{l=n+1}^{+\infty} M^{-(l-n)} P^{u} F_{l}  \tag{4.49}\\
& =\sum_{k=1}^{+\infty} M^{-k} P^{u} F_{k+n} \tag{4.50}
\end{align*}
$$

Then

$$
\begin{equation*}
\left\|J_{n}\right\|_{\mathbb{C}^{4}} \leq \sum_{k=1}^{+\infty}\left\|M^{-k} P^{u} F_{k+n}\right\|_{\mathbb{C}^{4}} \leq C \sum_{k=1}^{+\infty} \rho_{3}^{-k}\left\|F_{k+n}\right\|_{\mathbb{C}^{4}} \tag{4.51}
\end{equation*}
$$

with $\rho_{3}>1, C>0$, using the hyperbolicity assumption. Therefore

$$
\begin{equation*}
\left\|J_{n}\right\|_{\mathbb{C}^{4}} \leq C \sum_{-\infty}^{l=-1} \rho_{3}^{l}\left\|F_{n-l}\right\|_{\mathbb{C}^{4}}, \quad \forall n \geq 1 \tag{4.52}
\end{equation*}
$$

Letting $\alpha_{l}=\rho_{3}^{-|l|}$ if $l \leq-1, \alpha_{l}=0$ if $l \geq 0$, and $\beta_{m}=\left\|F_{m}\right\|_{\mathbb{C}^{4}}$ if $m \geq 1, \beta_{m}=0$ if $m \leq 0$, we then have that $\tilde{\zeta}_{n}$, given by

$$
\begin{equation*}
\tilde{\zeta}_{n}=C \sum_{l \in \mathbb{Z}} \alpha_{l} \beta_{n-l} \tag{4.53}
\end{equation*}
$$

is well defined for all $n \in \mathbb{Z}$, and satisfies $\tilde{\zeta}_{n}=\left\|J_{n}\right\|_{\mathbb{C}^{4}}$. We then have

$$
\begin{equation*}
\sum_{n=1}^{+\infty}\left\|J_{n}\right\|_{\mathbb{C}^{4}}^{2} \leq\|\tilde{\zeta}\|_{l^{2}}^{2} \leq C\left\|\rho_{3}^{-|\cdot|}\right\|_{l^{1}}^{2}\|F\|_{X_{c}^{4}}^{2} \tag{4.54}
\end{equation*}
$$

by Young's inequality (3.12). The left hand side is then finite, as required.
By (4.43), the existence of the limit implies $b_{+} \in E_{+}^{u}$. Then $P^{u} y_{1}=b_{+}$implies that $y_{1}=b_{+}+v_{+}^{s}$, where $v_{+}^{s}$ is an arbitrary element of $E_{+}^{s}$.

Similar arguments apply for the backwards solution of (4.36), starting from $y_{1}$, and iterating backwards for $n \leq 1$. We show that the projection to $E_{-}^{u}$ leads to a part of the solution that belongs to $X_{c,-}^{4}$ (where $X_{c,-}$ consists of sequences in $X_{c}$ that vanish
for $n>0$ ), while projection to $E_{-}^{s}$ leads to a term that can only decay if $y_{1}$ satisfies $y_{1}=b_{-}+v_{-}^{u}$, with $b_{-} \in E_{-}^{s}$ unique, and $v_{-}^{u}$ an arbitrary element of $E_{-}^{u}$.

Thus the two restrictions on $y_{1}$ imply that it belongs to two affine subspaces $b_{+}+E_{+}^{s}$, $b_{-}+E_{-}^{u}$. By the hypothesis on $E_{+}^{s}, E_{-}^{u}$, their intersection exists and is unique.

It is also easy to check that the solution $[p, q]^{T}$ of (4.34) obtained by gluing the forward and backward solution of (4.35) belongs to $X_{c} \times X_{c}$ if and only if the forward and backward solution of (4.35) belong to $X_{c,+}^{4}, X_{c,-}^{4}$ respectively.

Proof of Lemma 4.9. We will assume that $M_{+}(\lambda)$ has a semisimple eigenvalue $r_{c}$ on the unit circle, the other case is treated similarly. Fix an integer $m>4$, and $N=2 m+1$. Let $\tilde{z}_{n}^{N}, n \in \mathbb{Z}$ satisfy $\tilde{z}_{n}^{N}=0$, for all $n \leq 0, n>N$, and $\tilde{z}_{n+1}=M_{+} \tilde{z}_{n}$, for all $1 \leq n \leq N$, with $\tilde{z}_{1}^{N} \in E_{+}\left(r_{c}\right)$, the corresponding invariant subspace of $y_{n+1}=M_{+} y_{n}$.

By (4.20), (4.21) we have that $\left[(L-\lambda) \tilde{z}^{N}\right]_{n}=0$ for all $n$, except $n=0,1, N$, and $N+1$. Also $y_{1} \in E_{+}\left(r_{c}\right)$ implies that $\left\|y_{n+1}\right\|_{\mathbb{C}^{4}}=\left\|y_{n}\right\|_{\mathbb{C}^{4}}, \forall n \geq 1$. Then

$$
\begin{equation*}
\left|q_{2}\right|^{2}+\left|p_{2}\right|^{2}=\left|q_{2 p}\right|^{2}+\left|p_{2 p}\right|^{2}, \quad\left|q_{1}\right|^{2}+\left|p_{1}\right|^{2}=\left|q_{2 p+1}\right|^{2}+\left|p_{2 p+1}\right|^{2}, \quad \forall p \geq 1 . \tag{4.55}
\end{equation*}
$$

Letting $\left|q_{1}\right|^{2}+\left|p_{1}\right|^{2}=R_{1}^{2},\left|q_{2}\right|^{2}+\left|p_{2}\right|^{2}=R_{2}^{2}$, we calculate

$$
\begin{equation*}
\left\|\tilde{z}^{N}\right\|_{Y_{c}}^{2}=(m+1) R_{1}^{2}+m R_{2}^{2}, \tag{4.56}
\end{equation*}
$$

and

$$
\begin{equation*}
\left|\left[(L-\lambda) \tilde{z}^{N}\right]_{N}\right| \leq 3\left(|\omega+2|^{2}+|\lambda|^{2}\right) R_{1}^{2}+3 \delta^{2} R_{2}^{2}, \quad\left[(L-\lambda) \tilde{z}^{N}\right]_{N+1}=\delta^{2} R_{1}^{2} \tag{4.57}
\end{equation*}
$$

furthermore

$$
\begin{equation*}
\left[(L-\lambda) \tilde{z}^{N}\right]_{0}=\left[(L-\lambda) \tilde{z}^{N}\right]_{N+1}, \quad\left[(L-\lambda) \tilde{z}^{N}\right]_{1}=\left[(L-\lambda) \tilde{z}^{N}\right]_{N} . \tag{4.58}
\end{equation*}
$$

Letting $\tilde{u}^{N}=\left\|\tilde{z}^{N}\right\|_{Y_{c}}^{-1} \tilde{z}^{N}$ we therefore have

$$
\begin{equation*}
\left\|(L-\lambda) \tilde{u}^{N}\right\|_{Y_{c}} \leq \frac{C}{N} \tag{4.59}
\end{equation*}
$$

with $C$ depending on $R_{1}, R_{2}, \delta, \omega$, and $\lambda$. We therefore have a sequence of $\tilde{u}^{N}, N>9$, that belongs to $\operatorname{Ker}^{\perp}(L-\lambda)$ and satisfies

$$
\begin{equation*}
\left\|\tilde{u}^{N}\right\|_{Y_{c}}=1, \quad \forall N>0, \quad \text { and } \quad\left\|(L-\lambda) \tilde{u}^{N}\right\|_{Y_{c}} \rightarrow 0 \quad \text { as } \quad N \rightarrow \infty . \tag{4.60}
\end{equation*}
$$

Then $L-\lambda$, restricted to $\operatorname{Ker}^{\perp}(L-\lambda)$, can not have a closed range, this would imply the existence of a bounded inverse, contradicting (4.60). Thus $L-\lambda$ can not be Fredholm.
4.3. Essential spectrum of linearization around shelf-type breather: nonlocal case

We now examine the linearization around the solutions of Theorem 3.3 for the nonlocal case. We only analyze the essential spectrum. As in the local case we will compute the essential spectrum of a simpler operator that is a compact perturbation of the operator we are interested in. This simpler operator is also nonlocal, and we do not have a ready analogy with dynamical systems. This difference also yields a weaker result for the essential spectra of nonlocal cases, see Theorem 4.13.

To simplify the notation we will consider solutions obtained by continuation of the $\delta=0$ (finite $\kappa$ ) breather

$$
\widetilde{A}_{n}=\left\{\begin{array}{cc}
\alpha, & \text { if } n<0  \tag{4.61}\\
\frac{\alpha}{\sqrt{1-\rho}}, & \text { if } n=0 \\
0, & \text { if } n>0,
\end{array}\right.
$$

with $\alpha=\sqrt{-\frac{\omega}{2 \gamma}}, \rho=e^{-\kappa}$. This is a special case of the configuration of (3.9). For $\omega, \kappa$ fixed, let $B_{n}, n \in \mathbb{Z}$, be a solution of Theorem 3.3 and let

$$
\begin{equation*}
A_{n}=\widetilde{A}_{n}+B_{n}, \quad n \in \mathbb{Z} \tag{4.62}
\end{equation*}
$$

with $\widetilde{A}$ as in (4.61).
Let $\mathcal{L}_{\kappa, \delta}=J \mathcal{H}$, with $J \mathcal{H}$ as in (4.4), $L_{-}, L_{+}$as in (4.5), $\mathcal{A}, \mathcal{M}$ as in (4.6), (4.7), with $A_{n}$ as in (4.62). Operator $\mathcal{L}_{\kappa, \delta}=J \mathcal{H}$ is the linearization around the shelf-type solution of Theorem 3.3.

To calculate the essential spectrum of $\mathcal{L}_{\kappa, \delta}$ we will use its perturbation $\overline{\mathcal{L}}_{2}$, defined by $\overline{\mathcal{L}}_{2}=J \mathcal{H}$, where $J \mathcal{H}$ is as in (4.4) with $L_{-}, L_{+}$as in (4.5), and $\mathcal{A}, \mathcal{M}$ given by

$$
\begin{equation*}
\mathcal{A}(n, k)=\overline{\mathcal{A}}_{2}(n, k)=\tanh \frac{\kappa}{2} A_{n}^{2} \delta_{n, k}, \quad \mathcal{M}(n, k)=\overline{\mathcal{M}}_{2}(n, k)=\tanh \frac{\kappa}{2} A_{n} A_{m} e^{-\kappa|m-n|} \delta_{n, k}, \tag{4.63}
\end{equation*}
$$

for all $k, n \in \mathbb{Z}$, where $A_{n}, n \in \mathbb{Z}$, is as in (4.17), i.e. $A_{n}$ is the shelf solution of the local problem with $\delta=0$. We also define the operators $\overline{\mathcal{A}}_{1}, \overline{\mathcal{M}}_{1}$ by

$$
\begin{equation*}
\overline{\mathcal{A}}_{1}(n, k)=\tanh \frac{\kappa}{2} \widetilde{A}_{n}^{2} \delta_{n, k}, \quad \overline{\mathcal{M}}_{1}(n, k)=\tanh \frac{\kappa}{2} \widetilde{A}_{n} \widetilde{A}_{m} e^{-\kappa|m-n|} \delta_{n, k} \tag{4.64}
\end{equation*}
$$

for all $k, n \in \mathbb{Z}$, where $\widetilde{A}_{n}$ is as in (4.61). We easily see that $\overline{\mathcal{A}}_{1}-\overline{\mathcal{A}}_{2}, \overline{\mathcal{M}}_{1}-\overline{\mathcal{M}}_{2}$ are compact. For instance,

$$
\begin{equation*}
\left(\overline{\mathcal{M}}_{1}-\overline{\mathcal{M}}_{2}\right)(0, n)=\left(\overline{\mathcal{M}}_{1}-\overline{\mathcal{M}}\right)(n, 0)=\tanh \frac{\kappa}{2} \alpha^{2}\left((1-\rho)^{-1}-1\right) e^{-\kappa|n-1|}, \quad \forall n \leq 0 \tag{4.65}
\end{equation*}
$$

while all other entries of $\overline{\mathcal{M}}_{1}-\overline{\mathcal{M}}_{2}$ vanish. It easily follows that $\overline{\mathcal{M}}_{1}-\overline{\mathcal{M}}_{2}$ is HilbertSchmidt.
Lemma 4.12. The operator $\mathcal{L}_{\kappa, \delta}: Y \rightarrow Y$ is a compact perturbation of the operator $\overline{\mathcal{L}}_{2}: Y \rightarrow Y$.

Proof. To show that $\mathcal{L}_{\kappa, \delta}-\tilde{\mathcal{L}}$ is compact it is enough to show that $\mathcal{A}-\overline{\mathcal{A}}_{2}$, and $\mathcal{M}-\overline{\mathcal{M}}_{2}$ are compact.

The (squared) Hilbert-Schmidt norm of $T=\mathcal{A}-\overline{\mathcal{A}}_{2}$ is

$$
\begin{equation*}
\sum_{n \in \mathbb{Z}}\left|T \hat{e}_{n}\right|^{2}=\sum_{n \in \mathbb{Z}}\left|\tanh \frac{\kappa}{2} \sum_{m \in \mathbb{Z}} e^{\kappa|m-n|}\right| \widetilde{A}_{m}+\left.B_{m}\right|^{2}-\left.\bar{A}_{n}^{2}\right|^{2} \tag{4.66}
\end{equation*}
$$

where $\bar{A}_{n}=\frac{\alpha}{2}$, if $n \leq 0, \bar{A}_{n}=0$, if $n>0$. Using the equation (3.6) satisfied by $\widetilde{A}$, (4.66) becomes

$$
\begin{align*}
\sum_{n \in \mathbb{Z}}\left|T \hat{e}_{n}\right|^{2} & =\left.\left.\sum_{n \in \mathbb{Z}} 2\left|\tanh ^{2} \frac{\kappa}{2} 2 \sum_{m \in \mathbb{Z}} e^{\kappa|m-n|}\right| \widetilde{A}_{m} B_{m}\left|+\sum_{m \in \mathbb{Z}} e^{\kappa|m-n|}\right| B_{m}\right|^{2}\right|^{2}  \tag{4.67}\\
& \leq 2 \left\lvert\, \tanh ^{2} \frac{\kappa}{2}\left[2\left(\frac{\alpha}{1-\rho}\right)^{2}\left\|e^{-\kappa|\cdot|} *|B|\right\|_{l^{2}}^{2}+\left\|e^{-\kappa|\cdot|} *|B|^{2}\right\|_{l^{2}}^{2}\right]\right. \tag{4.68}
\end{align*}
$$

where $|B|_{n}=\left|B_{n}\right|,\left(|B|^{2}\right)_{n}=\left|B_{n}\right|^{2}, \forall n \in \mathbb{Z}$. We have also used $\left|\widetilde{A}_{n}\right| \leq \frac{\alpha}{1-\rho}, \forall n \in \mathbb{Z}$. The two terms of (4.68) are finite by Young's inequality, and the fact that $B \in l^{2}$. Thus $T$ is Hilbert-Schmidt. Since $\overline{\mathcal{M}}_{1}-\overline{\mathcal{M}}_{2}$ is a Hilbert-Schmidt, it is sufficient to show that
$\mathcal{M}-\overline{\mathcal{M}}_{1}$ is compact. The (squared) Hilbert-Schmidt norm of $S=\mathcal{M}-\overline{\mathcal{M}}_{1}$ is

$$
\begin{align*}
\sum_{n \in \mathbb{Z}}\left|S \hat{e}_{n}\right|^{2} & =\sum_{n \in \mathbb{Z}} \sum_{m \in \mathbb{Z}}\left|\tanh \frac{\kappa}{2} e^{-\kappa|m-n|}\left(\widetilde{A}_{m}+B_{m}\right)\left(\widetilde{A}_{n}+B_{n}\right)-\tanh \frac{\kappa}{2} \widetilde{A}_{m} \widetilde{A}_{n}\right|^{2}(  \tag{4.69}\\
& \leq 3 \tanh ^{2} \frac{\kappa}{2} \sum_{n \in \mathbb{Z}} \sum_{m \in \mathbb{Z}}\left(2\left|e^{-\kappa|m-n|} \widetilde{A}_{n} B_{m}\right|^{2}+\left|e^{-\kappa|m-n|} B_{n} B_{m}\right|^{2}\right)  \tag{4.70}\\
& \leq 3 \left\lvert\, \tanh ^{2} \frac{\kappa}{2}\left[2\left(\frac{\alpha}{1-\rho}\right)^{2}\left\|e^{-\kappa|\cdot|} *|B|\right\|_{l^{2}}^{2}+\left\|e^{-\kappa|\cdot|}\right\|_{l^{2}}^{2}\|B\|_{l^{2}}^{2}\right]\right. \tag{4.71}
\end{align*}
$$

using also $\left|\widetilde{A}_{n}\right| \leq \frac{\alpha}{1-\rho}, \forall n \in \mathbb{Z}$. The first term of (4.71) is finite by Young's inequality, and the fact that $B \in l^{2}$. Thus $S$ is also Hilbert-Schmidt.

Let $B_{1}$ be as in (4.30). Also for $\kappa>0, \rho=e^{-\kappa}$, define the set $B_{2, \kappa}$ by

$$
\begin{equation*}
\left.B_{2, \kappa}=\left\{-4 \delta \sin ^{2} \frac{k}{2}\left(2 \omega \frac{1-\rho^{2}}{1-2 \rho \cos k+\rho^{2}}+4 \delta \sin ^{2} \frac{k}{2}\right): \quad k \in \mathbb{R}\right)\right\} . \tag{4.72}
\end{equation*}
$$

Theorem 4.13. The essential spectrum of the operators $\overline{\mathcal{L}}_{2}$ and $\mathcal{L}_{\delta, \kappa}$ in $Y$ includes the set of $\lambda \in \mathbb{C}$ satisfying $\lambda^{2} \in B_{1} \cup B_{2, \kappa}$.

The two scenarios are suggested by the numerical results of [BCMP15], Figures 3, 4. These show calculations with finite shelf-type breathers with support of about $m=40$ sites (in the limit $\delta=0$ ), but we start to see the accumulation of eigenvalues in the regions indicated by $\lambda^{2} \in B_{1} \cup B_{2, \kappa}$.

Proof. By Lemma 4.12 it is enough to examine the essential spectrum of $\overline{\mathcal{L}}_{2}$. We compute approximate eigenfunctions of $\overline{\mathcal{L}}_{2}$, using the constant coefficient operators $\overline{\mathcal{L}}_{l}$, $\overline{\mathcal{L}}_{r}$ in $X \times X$ defined in the following way. Consider operators $A^{0}, M, h$ in $X$ with entries in the standard basis $M(n, m)=e^{-\kappa|n-m|}, A^{0}(n, m)=\frac{\omega}{2} \delta_{n, m}, n, m \in \mathbb{Z}$. Let $h(n, m)=\delta_{n, m}$ if $n, m \geq 0, h(n, m)=0$ otherwise. Then let

$$
\overline{\mathcal{L}}_{l}=\left[\begin{array}{cc}
0 & -\omega I-\delta \Delta+2 A^{0}  \tag{4.73}\\
\omega I+\delta \Delta-2 A^{0}-2 \omega M & 0
\end{array}\right], \quad \overline{\mathcal{L}}_{r}=\left[\begin{array}{cc}
0 & -\omega I-\delta \Delta \\
\omega I+\delta \Delta & 0
\end{array}\right] .
$$

In comparison, the definition of $\overline{\mathcal{L}}_{2}$ implies

$$
\overline{\mathcal{L}}_{2}=\left[\begin{array}{cc}
0 & -\omega I-\delta \Delta+2 h A^{0}  \tag{4.74}\\
\omega I+\delta \Delta-2 h A^{0}-2 \omega h M & 0
\end{array}\right] .
$$

Let $k \in \mathbb{R}$, and consider $z^{k} \in l^{\infty}(\mathbb{Z}, \mathbb{C})$ of the form

$$
\begin{equation*}
z_{n}^{k}=\left[a_{n} e^{i k n}+a_{n}^{*} e^{-i k n}, b_{n} e^{i k n}+b_{n}^{*} e^{-i k n}\right]^{T} . \tag{4.75}
\end{equation*}
$$

Then

$$
\begin{equation*}
\overline{\mathcal{L}}_{r} z^{k}=\lambda z^{k}, \quad \forall k \in \mathbb{Z} \tag{4.76}
\end{equation*}
$$

implies

$$
\begin{equation*}
\lambda^{2}=-\left(\omega-4 \delta \sin ^{2} \frac{k}{2}\right)^{2} \tag{4.77}
\end{equation*}
$$

Also,

$$
\begin{equation*}
\overline{\mathcal{L}}_{l} z^{k}=\lambda z^{k}, \quad \forall k \in \mathbb{Z} \tag{4.78}
\end{equation*}
$$

implies

$$
\begin{equation*}
\lambda^{2}=-4 \delta \sin ^{2} \frac{k}{2}\left(2 \omega \frac{1-\rho^{2}}{1-2 \rho \cos k+\rho^{2}}+4 \delta \sin ^{2} \frac{k}{2}\right) \tag{4.79}
\end{equation*}
$$

Let $k \in \mathbb{R}$, and $\lambda^{2}=\lambda^{2}(k)$ as in (4.77). For $\mu \geq 5$ integer, define $v_{\mu} \in X_{c} \times X_{c}$ by

$$
v_{\mu}(n)=\left\{\begin{array}{cl}
\frac{1}{\sqrt{2 \mu}} e^{i k n}[1,1]^{T}, & \text { if } \quad n \in\{1, \ldots, \mu\}  \tag{4.80}\\
0, & \text { if } n \in \mathbb{Z} \backslash\{1, \ldots, \mu\} .
\end{array}\right.
$$

Then $\left\|v_{\mu}\right\|_{Y_{c}}=1, \forall \mu \geq 5$. Note also that $A^{0} v_{\mu}=0$, and $M v_{\mu}=0, \forall \mu \geq 5$. Letting $\left[q_{\mu}, p_{\mu}\right]^{T}=v_{\mu}$, and using (4.74) we then have that

$$
\begin{equation*}
\left\|\left(\overline{\mathcal{L}}_{2}-\lambda\right) v_{m}\right\|_{Y_{c}}^{2}=\sum_{n \in \mathbb{Z}}\left(\left|\left[(-\omega I-\delta \Delta) p_{\mu}-\lambda q_{\mu}\right]_{n}\right|^{2}+\left|\left[(\omega I+\delta \Delta) q_{\mu}-\lambda p_{\mu}\right]_{n}\right|^{2}\right) \tag{4.81}
\end{equation*}
$$

By (4.76), (4.77) the only nonvanishing terms in the sum (4.81) are the ones corresponding to the four indices $n=0,1, \mu, \mu+1$. This holds for any $\mu \geq 5$. Each term is proportional to $\mu^{-1}$, and we easily see that

$$
\begin{equation*}
\left\|\overline{\mathcal{L}}_{2} v_{\mu}-\lambda v_{m}\right\|_{Y_{c}}^{2} \leq \frac{4}{\mu}(2|\omega|+7|\delta|)^{2} \rightarrow 0, \quad \text { as } \quad \mu \rightarrow \infty \tag{4.82}
\end{equation*}
$$

Thus $\lambda(k)$, with $\lambda^{2}$ as in (4.77), belongs to the essential spectrum of $\overline{\mathcal{L}}_{2}$. Varying $k$ over $\mathbb{R}$ we have that all $\lambda^{2} \in B_{1}$ belong to $\sigma_{\text {ess }}\left(\overline{\mathcal{L}}_{2}\right)$.

Let $k \in \mathbb{R}$, and $\lambda^{2}=\lambda^{2}(k)$ as in (4.79). For $\mu \geq 5$ integer, define $w_{\mu} \in X_{c} \times X_{c}$ by

$$
w_{\mu}(n)=\left\{\begin{array}{cc}
\frac{1}{\sqrt{2(\mu+1)}} e^{i k n}[1,1]^{T}, & \text { if } \quad n \in\{-\mu, \ldots, 0\}  \tag{4.83}\\
0, & \text { if } n \in \mathbb{Z} \backslash\{-\mu, \ldots, 0\}
\end{array}\right.
$$

We have $\left\|w_{\mu}\right\|_{Y_{c}}=1, \forall \mu \geq 5$.
We examine first the components $\left(\overline{\mathcal{L}}_{2} w_{\mu}-\lambda w_{\mu}\right)_{n}$, with $n>0$. We have $\left(A^{0} w_{\mu}\right)_{n}=0$, and $\left(M w_{\mu}\right)_{n}=0, \forall n>0$. Letting $\left[q_{\mu}, p_{\mu}\right]^{T}=w_{\mu},\left(\Delta q_{\mu}\right)_{n},\left(\Delta p_{\mu}\right)_{n}$ vanish for all $n>1$. We then compute

$$
\begin{align*}
& \sum_{n=1}^{\infty}\left(\left|\left[\left(-\omega I-\delta \Delta+2 h A^{0}\right) p_{\mu}-\lambda q_{\mu}\right]_{n}\right|^{2}+\right. \\
& \left.\quad+\left|\left[\left(\omega I+\delta \Delta-2 h A^{0}-2 \omega h M\right) q_{\mu}-\lambda p_{\mu}\right]_{n}\right|^{2}\right)=\frac{1}{\mu+1}, \tag{4.84}
\end{align*}
$$

for any $\mu \geq 5$.
Examining the components $\left(\tilde{\mathcal{L}} w_{\mu}-\lambda w_{\mu}\right)_{n}$ with $n \leq 0$, we have that by (4.78), (4.79)

$$
\begin{equation*}
\left[\left(-\omega I-\delta \Delta+2 h A^{0}\right) p_{\mu}-\lambda q_{\mu}\right]_{n}=0, \quad \forall n \in\{-\mu+1, \ldots,-1\}, \quad \text { and } \quad \forall n \leq-\mu-2 . \tag{4.85}
\end{equation*}
$$

For $n=0,-\mu,-\mu-1$ these terms are proportional to $(\mu+1)^{-1 / 2}$, and we compute

$$
\begin{align*}
\sum_{n=-\infty}^{0}\left|\left[\left(-\omega I-\delta \Delta+2 h A^{0}\right) p_{\mu}-\lambda q_{\mu}\right]_{n}\right|^{2} & \leq \frac{1}{\mu+1}\left[(3|\delta|+|\lambda|)^{2}+1\right]  \tag{4.86}\\
& \leq \frac{1}{\mu+1}\left[2|\delta|^{2}+4|\delta|\left(\frac{2|\omega|}{1-\rho}+4|\delta|\right)+1\right](4.87) \tag{4.87}
\end{align*}
$$

where $|\lambda|^{2}$ was bounded using (4.79). The above holds for any $\mu \geq 5$.
By $n \leq 0$ and the definition of $h A^{0}$ we also have

$$
\begin{equation*}
\left[\left(\omega I+\delta \Delta-2 h A^{0}-2 \omega h M\right) q_{\mu}-\lambda p_{\mu}\right]_{n}=\left[(\delta \Delta-2 \omega h M) q_{\mu}-\lambda p_{\mu}\right]_{n} \tag{4.88}
\end{equation*}
$$

Consider $n=-\mu-r, r \geq 2$. Then $\left[\Delta q_{\mu}\right]_{n}=\lambda p_{\mu}=0$. Also

$$
\begin{align*}
{\left[h M q_{\mu}\right]_{n} } & =\frac{1}{\sqrt{2(\mu+1)}} \sum_{m \in\{-\mu, \ldots, 0\}} \rho^{|-\mu-r-m|} e^{i k m}  \tag{4.89}\\
& =\frac{1}{\sqrt{2(\mu+1)}} \rho^{r} \rho^{\mu} \sum_{m \in\{-\mu, \ldots, 0\}} \rho^{m} e^{i k m} \\
& =\frac{1}{\sqrt{2(\mu+1)}} \rho^{r-1} \frac{\rho^{\mu+1}-e^{-i(\mu+1) k}}{1-\rho} \tag{4.90}
\end{align*}
$$

By (4.88), (4.90) we have

$$
\begin{align*}
\sum_{n=-\infty}^{-\mu-2}\left|\left[\left(\omega I+\delta \Delta-2 h A^{0}-2 \omega h M\right) q_{\mu}-\lambda p_{\mu}\right]_{n}\right|^{2} & \left.=\sum_{n=-\infty}^{-\mu-2} \mid 2 \omega h M q_{\mu}\right]\left._{n}\right|^{2}  \tag{4.91}\\
& =\sum_{r=2}^{\infty} \frac{2 \omega^{2} \rho^{2(r-1)}}{(\mu+1)} \frac{\left(\rho^{\mu+1}-e^{-i(\mu+1) k}\right)^{2}}{(1-\rho)^{2}} \\
& =\frac{2 \omega^{2}\left(\rho^{\mu+1}-e^{-i(\mu+1) k}\right)^{2} \rho^{2}}{(\mu+1)(1-\rho)^{2}\left(1-\rho^{2}\right)} \tag{4.92}
\end{align*}
$$

for any $\mu \geq 5$.
For $n=-\mu-1$, we can use (4.90) with $r=1$. We have

$$
\begin{equation*}
\left.[-2 \omega h M) q_{\mu}\right]_{-\mu-1}=-2 \omega \frac{1}{\sqrt{2(\mu+1)}} \frac{\rho^{\mu+1}-e^{-i(\mu+1) k}}{1-\rho} \tag{4.93}
\end{equation*}
$$

and therefore

$$
\begin{equation*}
\left|\left[\left(\omega I+\delta \Delta-2 h A^{0}-2 \omega h M\right) q_{\mu}-\lambda p_{\mu}\right]_{-\mu-1}\right|^{2}=\frac{1}{2(\mu+1)}\left(2|\omega|\left|\frac{\rho^{\mu+1}-e^{-i(\mu+1) k}}{1-\rho}\right|+|\delta|\right), \tag{4.94}
\end{equation*}
$$

for any $\mu \geq 5$.
We now consider $n \in\{-\mu+1, \ldots, 1\}$. To estimate

$$
\begin{equation*}
T_{n}=\left[\left(\omega I+\delta \Delta-2 h A^{0}-2 \omega h M\right) q_{\mu}-\lambda p_{\mu}\right]_{n}=\left[(\delta \Delta-2 \omega h M) q_{\mu}-\lambda p_{\mu}\right]_{n} \tag{4.95}
\end{equation*}
$$

we use the fact that $w_{\mu}$ satisfies (4.78), that is

$$
\begin{equation*}
\left[(\delta \Delta-2 \omega M) q_{\mu}-\lambda p_{\mu}\right]_{n}=0, \quad \forall n \in\{-\mu+1, \ldots, 1\} \tag{4.96}
\end{equation*}
$$

and (4.95) to write

$$
\begin{equation*}
T_{n}=\left[-2 \omega(h M-M) q_{\mu}\right]_{n}, \quad \forall n \in\{-\mu+1, \ldots, 1\} . \tag{4.97}
\end{equation*}
$$

By the definitions of $h$, and $M$ we have

$$
\begin{align*}
{\left[h M q_{\mu}\right]_{n} } & =(2 \mu+2)^{-1 / 2} \sum_{m \in\{-\mu, \ldots, 0\}} \rho^{|m-n|} e^{i k n}  \tag{4.98}\\
& =(2 \mu+2)^{-1 / 2}\left(\sum_{m \in\{n, \ldots, 0\}} \rho^{m-n} e^{i k n}+\sum_{m \in\{-\mu, \ldots, n-1\}} \rho^{n-m} e^{i k n}\right)  \tag{4.99}\\
& =(2 \mu+2)^{-1 / 2}\left(\sum_{m=n}^{0} \rho^{m-n} e^{i k(m-n)}+\sum_{m=-\mu}^{n-1} \rho^{n-m} e^{-i k(n-m)}\right) e^{i k n}(4  \tag{4.100}\\
& =(2 \mu+2)^{-1 / 2}\left(\sum_{\lambda=0}^{|n|} \rho^{\lambda} e^{i k \lambda}+\sum_{\lambda=1}^{\mu-|n|} \rho^{\lambda} e^{-i k \lambda}\right) e^{i k n} \tag{4.101}
\end{align*}
$$

and

$$
\begin{equation*}
\left[M q_{\mu}\right]_{n}=(2 \mu+2)^{-1 / 2}\left(\sum_{\lambda=0}^{\infty} \rho^{\lambda} e^{i k \lambda}+\sum_{\lambda=1}^{\infty} \rho^{\lambda} e^{-i k \lambda}\right) e^{i k n} \tag{4.102}
\end{equation*}
$$

By (4.97), (4.101), (4.102) we then have

$$
\begin{align*}
T_{n} & =\frac{2 \omega}{\sqrt{2(\mu+1)}}\left(\sum_{\lambda=|n|+1}^{\infty} \rho^{\lambda} e^{i k \lambda}+\sum_{\lambda=|\mu|-|n|+1}^{\infty} \rho^{\lambda} e^{-i k \lambda}\right) e^{i k n}  \tag{4.103}\\
& =\frac{2 \omega}{\sqrt{2(\mu+1)}}\left(\frac{\rho^{|n|+1} e^{i k(|n|+1)}}{1-\rho e^{i k}}+\frac{\rho^{\mu-|n|+1} e^{-i k(\mu-|n|+1)}}{1-\rho e^{-i k}}\right) e^{i k n} \tag{4.104}
\end{align*}
$$

Then

$$
\begin{equation*}
\sum_{n=-\mu+1}^{1}\left|T_{n}\right|^{2}=\frac{2 \omega}{\mu+1} \sum_{n=-\mu+1}^{1}\left(\frac{\rho^{2(|n|+1)}}{\left|1-\rho e^{i k}\right|^{2}}+\frac{\rho^{2(\mu-|n|+1)}}{\left|1-\rho e^{-i k}\right|^{2}}+2 \operatorname{Re} \frac{\rho^{\mu+2} e^{i k(\mu+2)}}{\left(1-\rho e^{i k}\right)\left(1-\rho e^{-i k}\right)}\right) \tag{4.105}
\end{equation*}
$$

The squared modulus terms are bounded by geometric series, e.g.

$$
\begin{equation*}
\left|1-\rho e^{i k}\right|^{-2} \sum_{n \in\{-\mu+1, \ldots, 1\}} \rho^{2(|n|+1)} \leq \frac{1}{(1-\rho)^{2}} \frac{1}{1-\rho^{2}} \tag{4.106}
\end{equation*}
$$

similarly for the second term. The sum of cross terms is bounded as

$$
\begin{equation*}
\sum_{n \in\{-\mu+1, \ldots, 1\}} 2 \operatorname{Re}\left(\rho^{\mu+2}\right) e^{i k(\mu+2)}\left(1-\rho e^{i k}\right)^{-1}\left(1-\rho e^{-i k}\right)^{-1} \leq \frac{2(\mu-1) \rho^{\mu+2}}{(1-\rho)^{2}} \tag{4.107}
\end{equation*}
$$

Then (4.105) yields

$$
\begin{equation*}
\sum_{n \in\{-\mu+1, \ldots, 1\}}\left|T_{n}\right|^{2} \leq \frac{2 \omega}{\mu+1}\left[\frac{2}{1-\rho^{2}}+2(\mu-1) \rho^{\mu+2} \frac{1}{(1-\rho)^{2}}\right] . \tag{4.108}
\end{equation*}
$$

This holds for all $\mu \geq 5$.

We finally consider the two sites $n=-\mu$, and $n=0$. At $n=-\mu$ we have

$$
\begin{align*}
{\left[h M q_{\mu}\right]_{-\mu} } & =\frac{1}{\sqrt{2(\mu+1)}} \sum_{m \in\{-\mu, \ldots, 0\}} \rho^{|-\mu-m|} e^{i k m}  \tag{4.109}\\
& =\frac{1}{\sqrt{2(\mu+1)}} e^{-i k \mu} \frac{1-\left(\rho e^{i k}\right)^{\mu+1}}{1-\rho e^{i k}} \tag{4.110}
\end{align*}
$$

therefore

$$
\begin{equation*}
\left|\left[(\delta \Delta-2 \omega h M) q_{\mu}-\lambda p_{\mu}\right]_{-\mu}\right|^{2} \leq \frac{1}{2(\mu+1)}\left(3|\delta|+|\lambda|+2|\omega| \frac{1}{1-\rho}\right) \tag{4.111}
\end{equation*}
$$

with $|\lambda|$ bounded in terms of $|\delta|,|\omega|$ as in (4.86). At $n=0,\left[h M q_{\mu}\right]_{0}$ is similarly

$$
\begin{equation*}
\left[h M q_{\mu}\right]_{0}=\frac{1}{\sqrt{2(\mu+1)}} \sum_{m \in\{-\mu, \ldots, 0\}} \rho^{|-m|} e^{i k m}=\frac{1}{\sqrt{2(\mu+1)}} \frac{1-\left(\rho e^{-i k}\right)^{\mu+1}}{1-\rho e^{-i k}} \tag{4.112}
\end{equation*}
$$

and

$$
\begin{equation*}
\left|\left[(\delta \Delta-2 \omega h M) q_{\mu}-\lambda p_{\mu}\right]_{0}\right|^{2} \leq \frac{1}{2(\mu+1)}\left(3|\delta|+|\lambda|+2|\omega| \frac{1}{1-\rho}\right) \tag{4.113}
\end{equation*}
$$

with $|\lambda|$ bounded in terms of $|\delta|,|\omega|$ as in (4.86). As before, (4.111), (4.113) are valid for any $\mu \geq 5$.

Collecting (4.84), (4.86), (4.94), (4.108), (4.111), (4.113) we see that there exists a $C$ that depends on $|\omega|,|\delta|$,and $\kappa$ for which

$$
\begin{equation*}
\left\|\overline{\mathcal{L}}_{2} \mu-\lambda w_{m}\right\|_{Y_{c}}^{2} \leq \frac{C}{\mu+1} \rightarrow 0, \quad \text { as } \quad \mu \rightarrow \infty \tag{4.114}
\end{equation*}
$$

thus $\lambda(k), \lambda^{2}$ as in (4.79), belongs to the essential spectrum of $\overline{\mathcal{L}}_{2}$. Varying over $k \in \mathbb{R}$ we have that all $\lambda^{2} \in B_{2, \kappa}$ belong to $\sigma_{\text {ess }}\left(\overline{\mathcal{L}}_{2}\right)$.

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