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## structure and DYnamics of noncommutative solitons: SPECTRAL THEORY AND DISPERSIVE ESTIMATES

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## ABSTRACT OF THE DISSERTATION

Structure and Dynamics of Noncommutative Solitons: Spectral Theory and Dispersive Estimates

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We consider the Schrödinger equation with a Hamiltonian given by a second order difference operator with nonconstant growing coefficients, on the half one dimensional lattice. This operator appeared first naturally in the construction and dynamics of noncommutative solitons in the context of noncommutative field theory. We prove pointwise in time decay estimates, with the optimal decay rate $t^{-1} \log ^{-2} t$ generically. We use a novel technique involving generating functions of orthogonal polynomials to achieve these estimates. We construct a ground state soliton for this equation and analyze its properties. In particular we arrive at $\ell^{\infty}$ and $\ell^{1}$ estimates as well as a quasi-exponential spatial decay rate. We completely determine the spectrum of the associated linearized Hamiltonian and prove the optimal decay rate of $t^{-1} \log ^{-2} t$ for the associated time decay estimate. These results are to appear in forthcoming papers.

## DEDICATION

This work is dedicated to my mother, without whom nothing would be possible.

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1 The dashed line is the graph of $f_{+}(a, b)$, the solid is of $f(a, b)$, and the dotted is of $f_{-}(a, b)$.

## Part 1

## Introduction and Background

The notion of noncommutative soliton arises when one considers the nonlinear Klein-Gordon equation (NLKG) for a field which is dependent on, for example, two "noncommutative coordinates", $x_{1}, x_{2}$, whose coordinate functions satisfy $\left[X_{1}, X_{2}\right]=$ $i \epsilon$. By going to a representation of the above canonical commutation relation, one can reduce the dynamics of the problem to an equation for the coefficients of an expansion in the Hilbert space representation of the above CCR, see e.g. [15]. By restricting to rotationally symmetric functions the nocommutative deformation of the Laplacian reduces to a second order finite difference operator, which is symmetric, and with variable coefficient growing like $n$, the lattice coordinate, at infinity. Therefore, this operator is unbounded, and in fact has continuous spectrum $[0, \infty)$. These preliminary analytical results, as well as additional numerical results, were obtained by Chen, Fröhlich, and Walcher. [5]. The dynamics and scattering of the (perturbed) soliton can then be inferred from the NLKG with such a discrete operator as the linear part. We will be interested in studying the dynamics of discrete NLKG and discrete NLS equations with these hamiltonians.

We follow the presentation of these methods in the manner of [5]. A soliton, or solitary wave, is a localized stationary solution to a nonlinear PDE, an example of which is the 2D real nonlinear Klein-Gordon equation: $-\partial_{t}^{2} \Phi=-\Delta \Phi+V(\Phi)$, where $\Phi: \mathbb{R}_{t} \times \mathbb{R}_{\left(x_{1}, x_{2}\right)}^{2} \rightarrow \mathbb{R}, V: \mathbb{R} \rightarrow \mathbb{R}$. A noncommutative soliton is a solitary wave solution to a PDE which has been deformed by enhancing the pointwise algebra of functions to one with a nontrivial spatial commutation relation. Here this is to say that $\left[X_{1}, X_{2}\right]=0 \mapsto\left[X_{1}, X_{2}\right]=i \epsilon$, where $X_{1}, X_{2}$ are the coordinate functions respectively for the Cartesian spatial coordinates $x_{1}, x_{2}$. By demanding that $\partial_{x} X=1$ one may arrive at $\partial_{x_{1}} \mapsto \frac{i}{\epsilon}\left[X_{2}, \cdot\right]$ and $\partial_{x_{2}} \mapsto \frac{i}{\epsilon}\left[X_{1}, \cdot\right]$. Finally, by considering formal power series expansions of functions on space one may deform the pointwise product $\Phi_{1} \Phi_{2}\left(x_{1}, x_{2}\right)$ to the Moyal star product

$$
\Phi_{1} \star \Phi_{2}\left(x_{1}, x_{2}\right)=e^{i(\epsilon / 2)\left(\partial_{a_{1}} \partial_{b_{2}}-\partial_{a_{2}} \partial_{b_{1}}\right)} \Phi_{1}\left(a_{1}, b_{1}\right) \Phi_{2}\left(a_{2}, b_{2}\right)\left\lfloor_{\left(a_{j}, b_{j}\right)=\left(x_{1}, x_{2}\right)} .\right.
$$

By analogy to $[Q, P]=i \epsilon \Rightarrow\left[A^{\dagger}, A\right]=1$ it is typical to consider the associated harmonic oscillator representation

$$
\begin{aligned}
& A:=(2 \epsilon)^{-1 / 2}\left(X_{1}+i X_{2}\right), \quad A^{\dagger}:=(2 \epsilon)^{-1 / 2}\left(X_{1}-i X_{2}\right), \\
& -\partial_{t}^{2} \Phi=-\Delta \Phi+V(\Phi) \quad \mapsto \quad-\partial_{t}^{2} \Phi=(2 / \epsilon)\left[A,\left[A^{\dagger}, \Phi\right]\right]+V(\Phi) .
\end{aligned}
$$

Let $\Pi_{n}$ be the projection onto the $n$-th excited state of the harmonic oscillator system: $\Pi_{n} v=v(n) \chi_{n}$. One may observe that

$$
\begin{aligned}
& \left(\chi_{n}, R^{2} \chi_{n}\right)=\left(\chi_{n},\left(X_{1}^{2}+X_{2}^{2}\right) \chi_{n}\right)=2 \epsilon(n-1 / 2), \\
& {\left[A,\left[A^{\dagger}, \Pi_{n}\right]\right]=\left\{\begin{array}{cc}
-(n+1) \Pi_{n+1}+(2 n+1) \Pi_{n}-n \Pi_{n-1} & , \quad n>0 \\
-\Pi_{1}+\Pi_{0} & n=0
\end{array}\right.}
\end{aligned}
$$

so that spherically symmetric functions are deformed into operators which are diagonal in the harmonic oscillator basis. Therefore if one restricts the original PDE to spherically symmetric solutions then the noncommutative PDE becomes a finite difference equation on the half-lattice of excited state level indices of a harmonic oscillator system: $\Phi(t, r) \mapsto \Phi(t, n)$.

As an alternative to noncommutative deformation one may consider the following. If one performs the change of variables $r^{2}=2 \epsilon n_{*}$, then the radial 2D laplacian takes the form $r^{-1} \partial_{r} r \partial_{r}=(2 / \epsilon) \partial_{n_{*}} n_{*} \partial_{n_{*}}$. Upon a careful choice of direct discretization one finds $(2 / \epsilon) \partial_{n_{*}} n_{*} \partial_{n_{*}} \mapsto(2 / \epsilon) D_{+} n D_{-}$, which has the same action on the $\chi_{n}$ as $-\Delta \mapsto(2 / \epsilon)\left[A,\left[A^{\dagger}, \cdot\right]\right]$ has on the $\Pi_{n}$, where $D_{+}, D_{-}$are respectively the forward and backward finite difference operators.

In this sense the resulting real radial 2D nonlinear Klein-Gordon equation can be considered as either a noncommutative version or a discrete version of the original commutative differential one. The result of the discretization procedure is perhaps a bit unexpected as the harmonic oscillator representation needn't look anything like the the original commutative spatial representation. Nevertheless, the original equation can be recovered with the limit $\epsilon \searrow 0$ as either the commutative limit
of a noncommutative system (like $\hbar \searrow 0$ ) or as the continuum limit of a lattice system. We will work with the finite difference representation exclusively. Henceforth $n \mapsto x \in \mathbb{Z}_{+}$, since the lattice will be our space, and we set $\epsilon=2$.

The principle of replacing the usual space with a so-called noncommutative space (or space-time) has found extensive use for model building in physics and in particular for allowing easier construction of localized solutions, see e.g. [2][22] for surveys. An example of the usefulness of this approach is that it may provide a robust procedure for circumventing classical nonexistence theorems for solitons, e.g. that of Derrick [8]. The NLKG variant of the equation we study here first appeared in the context of string theory and associated effective actions in the presence of background D-brane configurations, see e.g. [15]. We have decided to look in a completely different direction. The NLS variant and its solitons can in principle be materialized experimentally with optical devices, suitably etched, see e.g. [6]. Thus the dynamics of NLS with such solitons may offer new and potentially useful coherent states for optical devices. Furthermore, we believe the NLS solitons to have special properties, in particular asymptotic stability as opposed to the conjectured asymptotic metastability of the NLKG solitons conjectured in [5].

We will be following a procedure for the proof of asymptotic stability which has become standard within the study of nonlinear PDE. Crucial aspects of the theory and associated results were established by Buslaev and Perelman [3], Buslaev and Sulem [4], and Gang and Sigal [14]. Important elements of these methods are the dispersive estimates. Various such estimates have been found in the context of 1D lattice systems, for example see the work of A.I. Komech, E.A. Kopylova, and M. Kunze [20] and of I. Egorova, E. Kopylova, G. Teschl [12], as well as the continuum 2D problem to which our system bears many resemblances, see e.g. the work of E. A. Kopylova and A.I. Komech [21]. Extensive results have been found on the asymptotic stability on solitons of 1D nonlinear lattice Schrödinger equations by F. Palmero et al. [24] and P.G. Kevrekidis, D.E. Pelinovsky, and A. Stefanov [18]. Important aspects
of the application of these models to optical nonlinear waveguide arrays has been established by H.S. Eisenberg et al. [13].

This work is part of a series of forthcoming papers devoted to the construction, scattering, and asymptotic stability of these noncommutative solitons. The first three chapters are devoted to separate aspects of the problem in order of necessity. The fourth chapter addresses further work on this subject matter that we would like to pursue. The organization of this work is as follows.

In Chapter 1 we focus on a key estimate that is needed for scattering and stability, namely the decay in time of the solution, at the optimal rate. Fortunately, in the generic case, we find it is integrable, given by $t^{-1} \log ^{-2} t$. The proof of this result is rather direct, and employs the generating functions of the corresponding generalized eigenfunctions, to explicitly represent and estimate the resolvent of the hamiltonian at all energies. We also conclude the absence of positive eigenvalues and singular continuous spectrum.

Preliminary results for the scattering theory of the associated noncommutative waves and solitons were found by Durhuus and Gayral [9]. In particular they find local decay estimates for the associated noncommutative NLS. We utilize alternative methods and find local decay for both the free Schrödinger operator as well as a class of rank one perturbations thereof. An important element of this analysis is the study of the spectral properties of the free and perturbed Schrödinger operator. We extend the analysis of Chen, Frölich, and Walcher [5] and reproduce some of their results with alternative techniques.

In Chapter 2 we address the construction and properties of a family of ground state solitons. These stationary states satisfy a nonlinear eigenvalue equation, are positive, monotonically decaying and sharply peaked for large spectral parameter. The proof of this result follows directly from our spectral results in Chapter 1 by iteration for small data and root finding for large data. The existence and many properties of solutions for a similar nonlinear eigenvalue equation were found by Durhuus, Jonssen, and Nest
$[10][11]$. We utilize a simple power law nonlinearity for which their existence proofs do not apply. We additionally find estimates for the peak height, spatial decay rate, norm bounds, and parameter dependence.

In Chapter 3 we focus on deriving a decay rate estimate for the Hamiltonian which results from linearizing the original NLS around the soliton constructed in Chapter 1. We determine the full spectrum of this operator, which is the union of a multiplicity 2 null eigenvalue and a real absolutely continuous spectrum. This establishes a welldefined set of modulation equations [31] and points toward the asymptotic stability of the soliton.

In Chapter 4 we describe how the results from Chapter 3 can be applied to prove stability of the soliton we constructed in Chapter 2. The issue of asymptotic stability of NLS solitons has been sufficiently well-studied in such a broad context that the proof thereof is often considered as following straightforwardly from the appropriate spectral and decay estimates, of the kind found in Chapter 3. We sketch how the theory of modulation equations established by Soffer and Weinstein [31] can be used to prove asymptotic stability. Chen, Fröhlich, and Walcher conjectured that in the NLKG case the corresponding solitons are unstable but with exponential long decay: the so-called metastability property, see e.g. [5]. There is a great deal of evidence to suggest that this is in fact the case but a proof has yet to be provided. This will be the subject of future work.

## Part 2

## Main

## Notation

Let $\mathbb{Z}_{+}$and $\mathbb{R}_{+}$respectively be the nonnegative integers and nonnegative reals and $\mathscr{H}=\ell^{2}\left(\mathbb{Z}_{+}, \mathbb{C}\right)$ the Hilbert space of square integrable complex functions, e.g. $v: \mathbb{Z}_{+} \ni x \mapsto v(x) \in \mathbb{C}$, on the 1D half-lattice with inner product $(\cdot, \cdot)$, which is conjugate-linear in the first argument and linear in the second argument, and the associated norm $\|\cdot\|$, where $\|v\|=(v, v)^{1 / 2}, \forall v \in \mathscr{H}$. Where the distinction is clear from context $\|\cdot\| \equiv\|\cdot\|_{\text {op }}$ will also represent the norm for operators on $\mathscr{H}$ given by $\|A\|_{\text {op }}=\sup _{v \in \mathscr{H}}\|v\|^{-1}\|A v\|$, for all bounded $A$ on $\mathscr{H}$. Denote the lattice $\ell^{1}$ norm by $\|\cdot\|_{1}$ where $\|v\|_{1}=\sum_{x=0}^{\infty}|v(x)|, \forall v \in \ell^{1}\left(\mathbb{Z}_{+}, \mathbb{C}\right)$.

We denote by $\otimes$ the tensor product and by $z \mapsto \bar{z}$ complex conjugation for all $z \in \mathbb{C}$. We write $\mathscr{H}^{*}$ for the space of linear functionals on $\mathscr{H}$ : the dual space of $\mathscr{H}$. For every $v \in \mathscr{H}$ one has that $v^{*} \in \mathscr{H}^{*}$ is its dual satisfying $v^{*}(w)=(v, w)$ for all $v, w \in \mathscr{H}$. For every operator $A$ on $\mathscr{H}$ we take $\mathcal{D}(A)$ as standing for the domain of $A$. For each operator $A$ on $\mathscr{H}$ define $A^{*}$ on $\mathscr{H}^{*}$ to be its dual and $A^{\dagger}$ on $\mathscr{H}$ its adjoint such that $v^{*}(A w)=A^{*} v^{*}(w)=\left(A^{\dagger} v, w\right)$ for all $v \in \mathcal{D}\left(A^{\dagger}\right)$ and all $w \in \mathcal{D}(A)$. Let $\left\{\chi_{x}\right\}_{x=0}^{\infty}$ be the orthonormal set of vectors such that $\chi_{x}(x)=1$ and $\chi_{x_{1}}\left(x_{2}\right)=0$ for all $x_{2} \neq x_{1}$. We write $\Pi_{x}=\chi_{x} \otimes \chi_{x}^{*}$ for the orthogonal projection onto the space spanned by $\chi_{x}$.

We define $\mathscr{T}$ to be the topological vector space of all complex sequences on $\mathbb{Z}_{+}$ endowed with topology of pointwise convergence, $\mathcal{B}(\mathscr{H})$ to be the space of bounded linear operators on $\mathscr{H}$, and $\mathcal{L}(\mathscr{T})$ to be the space of linear operators on $\mathscr{T}$, endowed with the pointwise topology induced by that of $\mathscr{T}$. When an operator $A$ on $\mathscr{H}$ can be given by an explicit formula through $A\left(x_{1}, x_{2}\right)=\left(\chi_{x_{1}}, A \chi_{x_{2}}\right)<\infty$ for all $x_{1}, x_{2} \in \mathbb{Z}_{+}$ one may make the natural inclusion of $A$ into $\mathcal{L}(\mathscr{T})$, the image of which will also be
denoted by $A$. We consider $\mathscr{T}$ to be endowed with pointwise multiplication, i.e. the product $u v$ is specified by $(u v)(x)=u(x) v(x)$ for all $u, v \in \mathscr{T}$.

We represent the spectrum of each $A$ on $\mathscr{H}$ by $\sigma(A)$. We term each element $\lambda \in$ $\sigma(A)$ a spectral value. We write $\sigma_{\mathrm{d}}(A)$ for the discrete spectrum, $\sigma_{\mathrm{e}}(A)$ for the essential spectrum, $\sigma_{\mathrm{p}}(A)$ for the point spectrum, $\sigma_{\mathrm{ac}}(A)$ for the absolutely continuous spectrum, and $\sigma_{\mathrm{sc}}(A)$ for the singularly continuous spectrum. Should an operator satisfy the spectral theorem one may implement spectral projections. For each operator $A$, these will be written as $\Pi_{\mathrm{d}}^{A}$ and the like for each of the distinguished subsets of the spectral decomposition of A. Define $R_{.}^{A}: \rho(A) \rightarrow \mathcal{B}(\mathscr{H})$, the resolvent of $A$, to be specified by $R_{z}^{A}:=(A-z)^{-1}$, where $\rho(A):=\mathbb{C} \backslash \sigma(A)$ is the resolvent set of $A$ and where by abuse of notation $z I \equiv z \in \mathcal{B}(\mathscr{H})$ here.

Allow an eigenvector of $A$ to be a vector $v \in \mathscr{H}$ for which $A v=\lambda v$ for some $\lambda \in \mathbb{C}$. Should $A$ admit inclusion into $\mathcal{L}(\mathscr{T})$, we define a generalized eigenvector of $A$ be a vector $\phi \in \mathscr{T} \backslash \mathscr{H}$ which satisfies $A \phi=\lambda \phi$ for some $\lambda \in \mathbb{C}$ such that $\phi(x)$ is polynomially bounded, which is to say that there exists a $p \geq 0$ such that $\lim _{x \not \lambda_{\infty}}(x+1)^{-p} \phi(x)=0$. We define a spectral vector of $A$ to be a vector which is either an eigenvector or generalized eigenvector of $A$. We define the subspace of spectral vectors associated to the set $\Sigma \subseteq \sigma(A)$ to be the spectral space over $\Sigma$.

We write $\partial_{z} \equiv \frac{\partial}{\partial z}$ and $\mathrm{d}_{z} \equiv \frac{\mathrm{~d}}{\mathrm{~d} z}$ respectively for formal partial and total derivative operators with respect to a parameter $z \in \mathbb{R}, \mathbb{C}$.

## CHAPTER 1

## The linear scalar problem

## 1. Results

Definition 1.1. Define $L_{0}$ to be the operator on $\mathscr{H}$ with action

$$
L_{0} v(x)=\left\{\begin{array}{cc}
-(x+1) v(x+1)+(2 x+1) v(x)-x v(x-1) & , \quad x>0  \tag{1}\\
-v(1)+v(0) & , \quad x=0
\end{array}\right.
$$

and domain $\mathcal{D}\left(L_{0}\right):=\{v \in \mathscr{H} \mid\|M v\|<\infty\}$, where $M$ is the multiplication operator with action $M v(x)=x v(x) \forall v \in \mathscr{T}$.

Consider the linear Schrödinger equation

$$
\begin{equation*}
i \partial_{t} u=L_{0} u+V u \tag{2}
\end{equation*}
$$

where $u: \mathbb{R}_{t} \times \mathbb{Z}_{+} \rightarrow \mathbb{C}$ and $V$ is a potential (energy) multiplication operator on $\mathscr{H}$. To find solutions to Equation (2) it is sufficient to analyze the spectral measure of $L_{0}+V$.

Definition 1.2. For $A$ an essentially self adjoint operator on $\mathscr{H}$, such that each $\lambda \in \sigma(A)$ has multiplicity 1 , we denote the set of spectral vectors of $A$ by $\left\{\phi_{\lambda}^{A}\right\}_{\lambda \in \sigma(A)}$. Let a choice of normalization vector for $A$ be a fixed $v_{\mathrm{N}} \in \mathcal{H}$ such that $\left(v_{\mathrm{N}}, \phi_{\lambda}^{A}\right)=1$ for all $\lambda \in \sigma(A)$. Furthermore let the spectral integral weight $w_{\lambda}^{A}$ be given by $w_{\lambda}^{A}:=$ $\left(v_{\mathrm{N}}, \delta_{\lambda}^{A} v_{\mathrm{N}}\right)$ for all $\lambda \in \sigma(A)$. We define the resolvent vector $\psi^{A}: \rho(A) \rightarrow \mathscr{H}$ by $\psi_{z}^{A}:=$ $R_{z}^{A} v_{\mathrm{N}}$ and the auxilliary resolvent vector $\xi^{A}: \rho(A) \rightarrow \mathscr{H}$ by $\xi_{z}^{A}:=\psi_{z}^{A}-v_{\mathrm{N}}^{*}\left(\psi_{z}^{A}\right) \phi_{z}^{A}$, where $\phi_{z}^{A} \in \mathscr{T}$ is the analytic continuation $\sigma(A) \ni \lambda \rightarrow z \in \mathbb{C}$ of $\phi_{\lambda}^{A} \in \mathscr{T}$. Let the resolvent function $f^{A}: \rho(A) \rightarrow \mathbb{C}$ be specified by $\left(v_{\mathrm{N}}, R_{z}^{A} v_{\mathrm{N}}\right)$.

The above permits the useful representation $\psi_{z}^{A}=f_{z}^{A} \phi_{z}^{A}+\xi_{z}^{A} . \phi_{z}^{A}$ and $\psi_{z}^{A}$ are connected to the Stieltjes transform theory of orthogonal polynomials, in that $\phi_{z}^{A}$ is termed a primary polynomial and $\psi_{z}^{A}$ the corresponding secondary polynomial [33].

Proposition 1. The operator $L_{0}$ has the following properties.
(1) $L_{0}$ is essentially self-adjoint.
(2) The spectral space over each $\lambda \in \sigma\left(L_{0}\right)$ has multiplicity 1 .
(3) The spectrum of $L_{0}$ is absolutely continuous, $\sigma\left(L_{0}\right)=\sigma_{\mathrm{ac}}\left(L_{0}\right)=[0, \infty)$, and for choice of normalization $\phi_{\lambda}^{L_{0}}(0)=1$, its generalized eigenfunctions are the Laguerre polynomials $\phi_{\lambda}^{L_{0}}=\sum_{k=0}^{x} \frac{(-\lambda)^{k}}{k!}\binom{x}{k}$.

Chen, Frohlich, and Walcher determined the above properties for $L_{0}$ in [5] via methods which are different from ours with the apparent exception of a proof of essential selfadjointness.

Definition 1.3. Define $L$ to be the operator on $\mathscr{H}$ with domain $\mathcal{D}(L)=\mathcal{D}\left(L_{0}\right)$ and specified by $L:=L_{0}-q \Pi_{0}$ where $q \geq 0$ is a fixed constant.

Theorem 1. Let $\phi_{\lambda}^{L}, \lambda \in \sigma(L)$, denote spectral vectors of $L$ chosen to satisfy the normalization condition $\left(\chi_{0}, \phi_{\lambda}^{L}\right)=\phi_{\lambda}^{L}(0)=1, \forall \lambda \in \sigma(L)$. L has the following properties.
(1) $\sigma_{\mathrm{d}}(L)=\sigma_{\mathrm{p}}(L)=\left\{\lambda_{0}\right\}$, where $\lambda_{0}<0$ uniquely satisfies $1=q \psi_{\lambda_{0}}^{L_{0}}(0)$ and the unique eigenfunction over $\lambda_{0}$ is $\psi_{\lambda_{0}}^{L}=q \psi_{\lambda_{0}}^{L_{0}}$.
(2) $\sigma_{\mathrm{e}}(L)=\sigma_{\mathrm{ac}}(L)=[0, \infty)$.
(3) $\mathrm{d} \mu^{L}(\lambda) \Pi_{\mathrm{e}}^{L_{0}}=w_{\lambda}^{L} \phi_{\lambda}^{L} \otimes \phi_{\lambda}^{L, *} \mathrm{~d} \lambda \Pi_{\mathrm{e}}^{L_{0}}$, where $w_{\lambda}^{L}=\left\{\left[1-q e^{-\lambda} \mathcal{P} E_{1}(-\lambda)\right]^{2}+\right.$ $\left.\left[\pi q e^{-\lambda}\right]^{2}\right\}^{-1} e^{-\lambda}, \mathrm{d} \lambda$ is the Lebesgue measure on $[0, \infty)$, and $\mathcal{P} E_{1}(-\lambda)=$ $-\operatorname{Ei}(\lambda):=-\int_{-\lambda}^{\infty} \mathrm{d} u u^{-1} e^{-u}, \lambda>0$, is the exponential integral. The generalized eigenfunctions of $L$ are given by $\phi_{\lambda}^{L}=\phi_{\lambda}^{L_{0}}+q \xi_{\lambda}^{L_{0}}, \lambda \in \sigma_{\mathrm{ac}}(L)$.

We are ultimately interested in studying the solutions of a nonlinear equation so it is important to acquire decay estimates for dispersive "scattering states".

DEFINITION 1.4. Let $W_{\kappa, \tau}$ be the multiplication operator weight specified by $W_{\kappa, \tau} v(x)=(x+\kappa)^{\tau} v(x), \forall v \in \mathscr{T}$, where $0<\kappa \in \mathbb{R}, \tau \in \mathbb{R}$.

Theorem 2. For all $-3 \geq \tau \in \mathbb{R}, t>0, v \in \ell^{1}$, there exists a constant $c>0$ and a $1<\kappa \in \mathbb{R}$ such that

$$
\begin{equation*}
\left\|W_{\kappa, \tau} e^{-i t L_{0}} W_{\kappa, \tau} v\right\|_{\infty}<c t^{-1}\|v\|_{1} \tag{3}
\end{equation*}
$$

Theorem 3. For all $-3 \geq \tau \in \mathbb{R}, v \in \ell^{1}$, there exists a $1<\kappa \in \mathbb{R}$ such that

$$
\begin{equation*}
\left\|W_{\kappa, \tau} e^{-i t L} \Pi_{\mathrm{e}}^{L} W_{\kappa, \tau} v\right\|_{\infty}=\mathcal{O}\left(t^{-1} \log ^{-2} t\right), \quad t \nearrow \infty \tag{4}
\end{equation*}
$$

Our proof of these estimates will rely heavily on the generating functions of the generalized eigenvalues. This approach draws upon known properties of certain special functions. Hereby the problem of sequences on a lattice will be transformed into a problem of analytic functions in the complex plane.

## 2. Spectral Properties of $L_{0}$

Lemma 2.1. Any vector, $v$, the set of whose components, $\{v(x)\}_{x=0}^{\infty}$, have finitely many nonzero elements is a semi-analytic vector for $L_{0}$, which is to say that $\left\|L_{0}^{k} v\right\| \leq$ $c_{v}(2 k)$ ! where $c_{v}$ depends on $v$ alone.

Proof. Define $x_{v}:=\sup _{x}\{x: v(x) \neq 0\}$.

$$
\begin{align*}
\left\|L_{0} v\right\|_{2}^{2}= & \sum_{x=0}^{\infty}|-(x+1) v(x+1)+(2 x+1) v(x)-x v(x-1)|^{2}  \tag{5}\\
\leq & \sum_{x=0}^{\infty}[(x+1)|v(x+1)|+(2 x+1)|v(x)|+x|v(x-1)|]^{2}  \tag{6}\\
\leq & \sum_{x=0}^{\infty}\left\{\left[\left(x_{v}-1\right)+1\right]|v(x+1)|+\left(2 x_{v}+1\right)|v(x)|+\left(x_{v}+1\right)|v(x-1)|\right\}^{2} \\
= & \sum_{x=0}^{\infty}\left\{\left[x_{v}|v(x+1)|\right]^{2}+2 x_{v}\left(2 x_{v}+1\right)|v(x+1)||v(x)|\right. \\
& +2 x_{v}\left(x_{v}+1\right)|v(x) \| v(x-1)|+\left[\left(2 x_{v}+1\right)|v(x)|\right]^{2} \\
& \left.\quad+2\left(2 x_{v}+1\right)\left(x_{v}+1\right)|v(x)||v(x-1)|+\left[\left(x_{v}+1\right)|v(x-1)|\right]^{2}\right\} \\
\leq & \sum_{x=0}^{\infty}\left\{\left[x_{v}|v(x)|\right]^{2}+2 x_{v}\left(2 x_{v}+1\right)|v(x+1)||v(x)|\right. \\
& +2 x_{v}\left(x_{v}+1\right)|v(x+1)||v(x)|+\left[\left(2 x_{v}+1\right)|v(x)|\right]^{2} \\
& \left.+2\left(2 x_{v}+1\right)\left(x_{v}+1\right)|v(x+1)||v(x)|+\left[\left(x_{v}+1\right)|v(x)|\right]^{2}\right\} \\
\leq & 16\left(x_{v}+1\right)^{2}\|v\|_{1}^{2}
\end{align*}
$$

$$
\left\|L_{0} v\right\|_{1}=\sum_{x=0}^{\infty}|-(x+1) v(x+1)+(2 x+1) v(x)-x v(x-1)|
$$

$$
\leq \sum_{x=0}^{\infty}|(x+1) v(x+1)+(2 x+1) v(x)+x v(x-1)|
$$

$$
\leq \sum_{x=0}^{\infty}[(x+1)|v(x+1)|+(2 x+1)|v(x)|+x|v(x-1)|]
$$

$$
\begin{align*}
& \leq \sum_{x=0}^{\infty}\left\{\left[\left(x_{v}-1\right)+1\right]|v(x+1)|+\left(2 x_{v}+1\right)|v(x)|+\left(x_{v}+1\right)|v(x-1)|\right\}  \tag{18}\\
& \leq \sum_{x=0}^{\infty}\left[x_{v}|v(x+1)|+\left(2 x_{v}+1\right)|v(x)|+\left(x_{v}+1\right)|v(x-1)|\right]  \tag{19}\\
& \leq 4\left(x_{v}+1\right)| | v \|_{1} . \tag{20}
\end{align*}
$$

Let $a_{v}:=4\left(x_{v}+1\right)$. We have then that $\left\|L_{0} v\right\|_{1}^{2},\left\|L_{0} v\right\|_{2}^{2} \leq a_{v}\|v\|_{1}$. One may observe that $x_{L_{0} v}=x_{v}+1 \Rightarrow a_{L_{0} v}=a_{v}+4$. One then has that

$$
\begin{align*}
\left\|L_{0}^{k} v\right\|_{2} & \leq a_{L_{0}^{k-1} v}\left\|L_{0}^{k-1} v\right\|_{1} \leq a_{L_{0}^{k-1} v} a_{L_{0}^{k-2} v}\left\|L_{0}^{k-2} v\right\|_{1} \leq \ldots  \tag{21}\\
& \leq \prod_{j=1}^{k} a_{L_{0}^{k-j} v}\|v\|_{1}=4^{k} \prod_{j=1}^{k}\left(j+x_{v}\right)\|v\|_{1}=4^{k}\left(x_{v}!\right)^{-1}\left(k+x_{v}\right)!\|v\|_{1} . \tag{22}
\end{align*}
$$

Since $4^{k}\left(x_{v}!\right)^{-1}\left(k+x_{v}\right)!\|v\|_{1}$ is monotonically increasing in $k \in \mathbb{Z}_{+}$we may, without loss of generality, take that $k>x_{v}$ and $k>4$ to bound this expression. One may show through the monotonicity in $k$ of $\binom{x}{k}$ for $0 \leq k \leq\lfloor x / 2\rfloor$ that $\binom{x}{k} \leq(\lfloor x / 2\rfloor!)^{-2}(x!)$, where $\lfloor a\rfloor=\sup _{a \geq n \in \mathbb{Z}} n$, for all $a \in \mathbb{R}$, is the floor function. One then has

$$
\begin{align*}
4^{k}\left(x_{v}!\right)^{-1}\left(k+x_{v}\right)!\|v\|_{1} & =4^{k} k!\binom{k+x_{v}}{k}\|v\|_{1} \leq 4^{k} k!\binom{2 k}{k}\|v\|_{1}  \tag{23}\\
& =4^{k}(k!)^{-1}(2 k)!\|v\|_{1} \leq 4^{4}(4!)^{-1}(2 k)!\|v\|_{1}<c_{v}(2 k)! \tag{24}
\end{align*}
$$

where $c_{v}=11\|v\|_{1}$.

Proof of Proposition (1) Part (1). The set of vectors, $v$, with finitely many nonzero components is dense in $\mathscr{H}$. This dense set is semi-analytic for $L_{0}$. By the extension of Nelson's analytic vector theorem to semi-analytic vectors, see e.g. [28], it is therefore the case that $L_{0}$ is essentially self-adjoint.

Proof of Proposition (1) Part (2). One has that $L_{0} v(x)=\lambda v(x), v \in \mathscr{T}$, specifies a countable family of coupled elementary algebraic equations. A unique solution may be found for each $\lambda$ by specifying $v_{0}$ and solving inductively in increasing
$x \in \mathbb{Z}_{+}$. A choice of normalization will fix $v(0)$. Therefore each solution is, up to normalization, uniquely specified by $\lambda$.

Proof of Proposition (1) Part (3). $L_{0}$ is an essentially self-adjoint, second order, finite difference operator or Jacobi operator. It is well known that the theory of Jacobi operators is intimately connected with that of orthogonal polynomials. In particular spectral equations for operators extended to formal sequence spaces may be viewed as recursion formulas for families, indexed by lattice site, of orthogonal polynomials defined on the spectrum of the operator in question, see e.g. [27]. For $L_{0} \in \mathcal{L}(\mathscr{T})$ it is the case that $L_{0} v(x)=\lambda v(x)$ takes the form of the recursion formula for the Laguerre polynomials. By part (2) of the proposition these solutions are unique.

The Laguerre polynomials, $\phi_{\lambda}^{L_{0}}(x) \equiv \phi_{\lambda}(x)$, have known completeness and orthogonality relations whose roles will be reversed here:

$$
\begin{equation*}
\delta_{x_{1}, x_{2}}=\int_{0}^{\infty} \mathrm{d} \lambda e^{-\lambda} \phi_{\lambda}\left(x_{1}\right) \phi_{\lambda}\left(x_{2}\right), \quad \delta\left(\lambda_{1}-\lambda_{2}\right)=e^{-\left(\lambda_{1}+\lambda_{2}\right) / 2} \sum_{n=0}^{\infty} \phi_{\lambda_{1}}(x) \phi_{\lambda_{2}}(x) \tag{25}
\end{equation*}
$$

where here $\delta(\cdot)$ is Dirac's delta distribution supported on $\sigma\left(L_{0}\right)$. The RHS of these equations converge in the distributional sense respectively on $\ell^{2}\left(\mathbb{Z}_{+}\right)$and $L^{2}\left(\mathbb{R}_{+}\right)$. The former equation expresses components of the spectral measure of $L_{0}$ and in particular $w_{\lambda}^{L_{0}}=e^{-\lambda}$. In what follows we will often suppress the dependence of a quantity on $L_{0}$, e.g. $w_{\lambda}^{L_{0}} \equiv w_{\lambda}$. We will use $\chi_{0}$ a normalization vector for $L_{0}$.

## 3. Spectral Properties of $L$

Consider that $A_{0}$ is an essentially self adjoint operator on $\mathscr{H}$ and has spectrum $\sigma\left(A_{0}\right)=\sigma_{\mathrm{ac}}\left(A_{0}\right)=[0, \infty)$. Furthermore consider that $A$ on $\mathcal{H}$ has domain $\mathcal{D}(A)=$ $\mathcal{D}\left(A_{0}\right)$ and is specified by $A:=A_{0}-q \Pi$, where $q \in \mathbb{C}$ is a fixed constant and $\Pi$ is a rank-1 orthogonal projection operator. There exists a $v_{\mathrm{P}} \in \mathscr{H}$ of unit norm for which $\Pi=v_{\mathrm{P}} \otimes v_{\mathrm{P}}^{*}$. Define $f_{z}^{A}:=\left(v_{\mathrm{P}}, R_{z}^{A} v_{\mathrm{P}}\right)$. By Weyl's theorem, the perturbation $-q \Pi$ cannot change the essential spectrum of $A_{0}$ but it can introduce eigenvalues. Let $u \in \mathscr{H}$ satisfy $v_{\mathrm{P}}^{*}(u) \neq 0$, then

$$
\begin{equation*}
A u=z u \quad \Rightarrow \quad 1=q f_{z}^{A} \tag{26}
\end{equation*}
$$

$A$ will then have as many eigenvalues as $q f_{z}^{A}-1$ has zeroes. The corresponding eigenfunctions are then given by

$$
\begin{equation*}
A u=\lambda u \quad \Rightarrow \quad u=q v_{\mathrm{P}}^{*}(u) R_{\lambda}^{A_{0}} v_{\mathrm{P}} . \tag{27}
\end{equation*}
$$

Definition 3.1. Let $\mathcal{T} \in \mathcal{L}(\mathscr{T})$ be the binomial transform, see e.g. [19], defined by

$$
\begin{equation*}
\mathcal{T} v(k)=\sum_{x=0}^{\infty} \mathcal{T}(k, x) v(x)=\sum_{x=0}^{\infty}(-1)^{k}\binom{k}{x} v(x), \quad \forall v \in \mathscr{T} . \tag{28}
\end{equation*}
$$

$\mathcal{T}$ is involutive in the sense that $\mathcal{T}^{2}=I$. One has that $\mathcal{T} v(0)=v(0)$ and the useful representation $\chi_{0}(x)=\sum_{k=0}^{\infty}(-1)^{x}\binom{x}{k}$. We take the conventions that $x!,\binom{x}{k}, \sum_{x=0}^{k} v(x)=0$ for $k, x<0$ and $k<x$ for all $v \in \mathscr{T}$.

Lemma 3.1. One may check by direct computation that

$$
\begin{equation*}
\mathcal{T} L_{0} v(k)=(k+1) \mathcal{T} v(k+1), \quad \forall v \in \mathscr{T} \tag{29}
\end{equation*}
$$

and may recover the usual definition of the Laguerre polynomials hereby.

Proof. Consider

$$
\begin{equation*}
L_{0} \phi_{\lambda}(x)=\lambda \phi_{\lambda}(x) \Rightarrow(k+1) \mathcal{T} \phi_{\lambda}(k+1)=\lambda \mathcal{T} \phi_{\lambda}(k) \tag{30}
\end{equation*}
$$

Choosing $\mathcal{T} \phi_{\lambda}(0)=\phi_{\lambda}(0)=1$ one has by induction that $\mathcal{T} \phi_{\lambda}(k)=\frac{\lambda^{k}}{k!}$. One may then apply the binomial transform again to arrive at

$$
\begin{equation*}
\phi_{\lambda}(x)=\sum_{k=0}^{x} \frac{(-\lambda)^{k}}{k!}\binom{x}{k} . \tag{31}
\end{equation*}
$$

Lemma 3.2. One has the representation

$$
\begin{equation*}
\psi_{z}(x)=e^{-z} \sum_{k=0}^{x}(-1)^{k}\binom{x}{k} E_{k+1}(-z) \tag{32}
\end{equation*}
$$

where

$$
\begin{equation*}
E_{p}(z):=\int_{1}^{\infty} \mathrm{d} t e^{-z t} t^{-p}, \quad p \in \mathbb{C}, z \in \mathbb{C} \backslash(-\infty, 0] \tag{33}
\end{equation*}
$$

are the generalized exponential integrals for which we take the principal branch with standard branch cut $\Sigma=(-\infty, 0]$.

Proof. Consider the $\left(L_{0}-z\right) \psi_{z}=\chi_{0}$, where $\psi_{z}, \chi_{0} \in \mathscr{T}, L_{0} \in \mathcal{L}(\mathscr{T})$. By binomial transform of this equation one finds

$$
\begin{equation*}
(k+1) \mathcal{T} \psi_{z}(k+1)=z \mathcal{T} \psi_{z}(k)+1 \tag{34}
\end{equation*}
$$

$\mathcal{T} \psi_{z}(k)=e^{-z} E_{k+1}(-z)$ satisfies this recursion formula.

Definition 3.2. For any single-valued or multi-valued function $f: \mathbb{C} \rightarrow \mathbb{C}$, an element of a set of linear functionals on some suitable Banach space with norm given through integration over $\lambda$, and with poles, branch points, and branch cuts found in
the subset $\Sigma \subseteq \mathbb{R}$ let $\mathcal{P} f: \Sigma \rightarrow \mathbb{C}$ be the principal part of $f$ defined by the weak limit

$$
\begin{equation*}
\mathcal{P} f(\lambda):=\frac{1}{2} \underset{\epsilon}{\mathrm{w}-\lim _{0}}[f(\lambda+i \epsilon)+f(\lambda-i \epsilon)], \quad \lambda \in \Sigma \tag{35}
\end{equation*}
$$

which converges in the distributional sense. We analogously define the $\delta$-part of $f$ to be

$$
\begin{equation*}
\delta f(\lambda):=\frac{1}{2 \pi i} \underset{\epsilon \backslash 0}{\mathrm{w}-\lim _{0}}[f(\lambda+i \epsilon)-f(\lambda-i \epsilon)], \quad \lambda \in \Sigma . \tag{36}
\end{equation*}
$$

We have kept vague the specification of the sense in which the above definitions converge weakly for the purposes of generality. The details of such convergence in our work will be clear from context. One may extend the domain of $\mathcal{P} f$ to the complex plane and produce a single valued function, which we will also denote $f$, through

$$
\mathcal{P} f(z):=\left\{\begin{array}{rr}
f(z), & z \in \mathbb{C} \backslash \Sigma  \tag{37}\\
\mathcal{P} f(z), & z \in \Sigma
\end{array}\right.
$$

One may observe that the analogous extension of $\delta f(\lambda)$ vanishes away from $\Sigma \subseteq \mathbb{R}$. This prescription extends to weak limits in $z \in \mathbb{C}$ of complex sequences $v_{z} \in \mathscr{T}$ whose components depend upon $z$.

The generalized exponential integrals have the convergent series expansion [34]

$$
\begin{align*}
E_{n+1}(z)=- & \frac{(-z)^{n}}{n!} \log (z)+\frac{e^{-z}}{n!} \sum_{k=1}^{n}(-z)^{k-1}(n-k)!  \tag{38}\\
& +\frac{e^{-z}(-z)^{n}}{n!} \sum_{k=0}^{\infty} \frac{z^{k}}{n!} \digamma(k+1) \tag{39}
\end{align*}
$$

where $\digamma(x):=\mathrm{d}_{x} \log \Gamma(x)$ is the digamma function. One may therefore observe that

$$
\begin{equation*}
\underset{\epsilon \searrow 0}{\mathrm{w}-\lim _{0}} E_{n+1}(-x \pm i \epsilon)=\mathcal{P} E_{n+1}(-x) \mp i \pi \frac{(x)^{n}}{n!}, \quad x>0 \tag{40}
\end{equation*}
$$

where for the sake of generality the limit is weak with respect to $L^{2}([a, \infty), \mathbb{C}), a>0$. In particular one often writes $\mathcal{P} E_{1}(-x)=-\operatorname{Ei}(x)$ where

$$
\begin{equation*}
\operatorname{Ei}(x):=-\int_{-x}^{\infty} \mathrm{d} u u^{-1} e^{-u}, \quad x>0 \tag{41}
\end{equation*}
$$

is the exponential integral.

Proof of Theorem (1) Part (1). Eigenvalues are be given by zeros of $q e^{-z} E_{1}(-z)-1$. $L$ is essentially self-adjoint so $\sigma(L) \subseteq \mathbb{R}$. Analytic continuation of $E_{1}(-z)$ to $z>0$ from above or below will result in the sum of a real function and an imaginary constant

$$
\begin{equation*}
\lim _{\epsilon \searrow 0} E_{1}(-x \pm i \epsilon)=-\operatorname{Ei}(x) \mp i \pi, \quad x>0 \tag{42}
\end{equation*}
$$

so there can be no positive eigenvalues. $q e^{-z} E_{1}(-z)$ diverges for $z \rightarrow 0$ so $z=0$ cannot be an eigenvalue. All eigenvalues must be negative. Let $z=-a<0$. It is the case that $e^{a} E_{1}(a)$ is monotonically decreasing for increasing $a \in(0, \infty]$

$$
\begin{equation*}
\mathrm{d}_{a}\left[e^{a} E_{1}(a)\right]=-\int_{0}^{\infty} \mathrm{d} x e^{-x}(x+a)^{-2}<0 \tag{43}
\end{equation*}
$$

where we have used manifest dominated convergence of the integral to pass the derivative through the integral. Furthermore since

$$
\begin{equation*}
\lim _{a \searrow 0} \int_{0}^{\infty} \mathrm{d} t e^{-t}(t+a)^{-1}=\infty, \quad \lim _{a \nearrow \infty} \int_{0}^{\infty} \mathrm{d} t e^{-t}(t+a)^{-1}=0 \tag{44}
\end{equation*}
$$

it follows that $e^{a} E_{1}(a)$ takes each on the interval $[0, \infty)$ exactly once, where we have used manifest uniform convergence of the integrand to pass the limit through the integral. Therefore $q e^{-z} E_{1}(-z)-1$ has exactly one root for each fixed $q>0$.

Proof of Theorem (1) Part (2). By the argument of the Proof of Theorem (1) Part (1), there can be no embedded eigenvalues. By Weyl's critereon the perturbation of $L_{0} \mapsto L$ leaves the essential spectrum unchanged. The argument for the proof of $\sigma\left(L_{0}\right)=\sigma_{\mathrm{ac}}\left(L_{0}\right)$ follows without change for the spectrum of $L$.

Definition 3.3. Let $A$ be an operator on $\mathscr{H}$ which is self-adjoint on its domain $\mathcal{D}(A)$ and $\lambda$ an element of the discrete spectrum of $A$. Define $\mathcal{P}_{\lambda}^{A} \equiv \mathcal{P}(A-\lambda)^{-1}$, $\lambda \in \sigma(A)$ to be the principal part of the resolvent of $A$ given by the strong limit

$$
\begin{equation*}
\mathcal{P}_{\lambda}^{A}:=\frac{1}{2} \underset{\epsilon}{\mathrm{~s}-\lim _{\star}}\left[R_{\lambda+i \epsilon}^{A}+R_{\lambda-i \epsilon}^{A}\right] . \tag{45}
\end{equation*}
$$

Denote by $\delta_{\lambda}^{A} \equiv \delta(A-\lambda) \equiv \Pi_{\lambda}^{A}, \lambda \in \sigma(A)$ the spectral projection defined by the strong limit

$$
\begin{equation*}
\delta_{\lambda}^{A}:=\frac{1}{2 \pi i} \mathrm{~s}_{\epsilon}^{\epsilon} \lim _{0}\left[R_{\lambda+i \epsilon}^{A}-R_{\lambda-i \epsilon}^{A}\right] . \tag{46}
\end{equation*}
$$

If $\lambda$ is instead an element of the essential spectrum of $A$ one has that $\mathcal{P}_{\lambda}^{A}, \delta_{\lambda}^{A}$ are defined by weak limits. If and only if the essential spectrum of $A$ is absolutely continuous it is the case that $\mathrm{d} \mu_{\mathrm{e}}^{A}(\lambda)=\delta_{\lambda}^{A} \mathrm{~d} \lambda$, where $\mathrm{d} \mu_{\mathrm{e}}^{A}(\lambda)$ is the essential spectral measure of $A$ and $\mathrm{d} \lambda$ is the Lebesgue measure on $\sigma_{\mathrm{e}}(A)$.

The above definition permits the useful representation $\delta_{\lambda}^{A}=w_{\lambda}^{A} \phi_{\lambda}^{A} \otimes \phi_{\lambda}^{A, *}$. One may typically observe through the spectral representation of $R_{z}^{A}$ that $\mathcal{P} \psi_{\lambda}^{A}=\mathcal{P}_{\lambda}^{A} v_{N}$ and that $\mathcal{P} \xi_{\lambda}^{A}=\mathcal{P} \psi_{\lambda}^{A}-v_{\mathrm{N}}^{*}\left(\mathcal{P} \psi_{\lambda}^{A}\right) \phi_{\lambda}^{A}=\xi_{\lambda}^{A}, \forall \lambda \in \sigma(A)$.

We recall the method of spectral shifts as applied to rank-1 perturbations, see e.g. [29], for operators of the form specified by the $A_{0}, \Pi, A=A_{0}-q \Pi$ considered above. Through the resolvent formula it follows that

$$
\begin{align*}
R_{z}^{A} & =R_{z}^{A_{0}}+R_{z}^{A_{0}} q \Pi R_{z}^{A} \quad \Rightarrow \quad \Pi R_{z}^{A}=\Pi R_{z}^{A_{0}}+f_{z}^{A_{0}} q \Pi R_{z}^{A}  \tag{47}\\
\Rightarrow \quad \Pi R_{z}^{A} & =\left(1-q f_{z}^{A_{0}}\right)^{-1} \Pi R_{z}^{A_{0}} \quad \Rightarrow \quad R_{z}^{A}=R_{z}^{A_{0}}+\left(1-q f_{z}^{A_{0}}\right)^{-1} R_{z}^{A_{0}} q \Pi R_{z}^{A_{0}} \tag{48}
\end{align*}
$$

For $A$ essentially self-adjoint one may apply the definitions of $\mathcal{P}_{\lambda}^{A}$ and $\delta_{\lambda}^{A}$ and find the corresponding shifts to $\mathcal{P}_{\lambda}^{A_{0}}$ and $\delta_{\lambda}^{A_{0}}$. For $\lambda \in \sigma(A)$ it follows that

$$
\begin{align*}
\mathcal{P}_{\lambda}^{A}= & \mathcal{P}_{\lambda}^{A_{0}}+g_{\lambda}^{A_{0}}\left[\left(1-q \mathcal{P} f_{\lambda}^{A_{0}}\right)\left(\mathcal{P}_{\lambda}^{A_{0}} q \Pi \mathcal{P}_{\lambda}^{A_{0}}-\pi^{2} \delta_{\lambda}^{A_{0}} q \Pi \delta_{\lambda}^{A_{0}}\right)\right.  \tag{49}\\
& \left.-\pi^{2} q \delta f_{\lambda}^{A_{0}}\left(\mathcal{P}_{\lambda}^{A_{0}} q \Pi \delta_{\lambda}^{A_{0}}+\delta_{\lambda}^{A_{0}} q \Pi \mathcal{P}_{\lambda}^{A_{0}}\right)\right] \tag{50}
\end{align*}
$$

$$
\begin{gather*}
\delta_{\lambda}^{A}=\delta_{\lambda}^{A_{0}}+g_{\lambda}^{A_{0}}\left[\left(1-q \mathcal{P} f_{\lambda}^{A_{0}}\right)\left(\mathcal{P}_{\lambda}^{A_{0}} q \Pi \delta_{\lambda}^{A_{0}}+\delta_{\lambda}^{A_{0}} q \Pi \mathcal{P}_{\lambda}^{A_{0}}\right)\right.  \tag{51}\\
\left.+q \delta f_{\lambda}^{A_{0}}\left(\mathcal{P}_{\lambda}^{A_{0}} q \Pi \mathcal{P}_{\lambda}^{A_{0}}-\pi^{2} \delta_{\lambda}^{A_{0}} q \Pi \delta_{\lambda}^{A_{0}}\right)\right], \tag{52}
\end{gather*}
$$

where

$$
\begin{equation*}
g_{\lambda}^{A_{0}}:=\left[\left(1-q \mathcal{P} f_{\lambda}^{A_{0}}\right)^{2}+\left(q \pi \delta f_{\lambda}^{A_{0}}\right)^{2}\right]^{-1} \tag{53}
\end{equation*}
$$

If $v_{\mathrm{P}}=v_{\mathrm{N}}$ is a normalizing vector for $A_{0}$ one may greatly simplify the expression of the shift of the spectral projection. Since one has $f_{z}^{A}=\left(v_{\mathrm{N}}, R_{z}^{A} v_{\mathrm{N}}\right)$ it follows that $\mathcal{P} f_{\lambda}^{A}=v_{\mathrm{N}}^{*}\left(\mathcal{P} \psi_{\lambda}^{A}\right)=\mathcal{P} v_{\mathrm{N}}^{*}\left(\psi_{\lambda}^{A}\right)$ and $\delta f_{\lambda}^{A}=w_{\lambda}^{A}$. One may find that

$$
\begin{align*}
R_{z}^{A} \alpha & =\psi_{z}^{A}=\left(1-q f_{z}^{A_{0}}\right)^{-1} \psi_{z}^{A_{0}} \Rightarrow f_{z}^{A}=\left(1-q f_{z}^{A_{0}}\right)^{-1} f_{z}^{A_{0}}  \tag{54}\\
\mathcal{P} f_{z}^{A} & =g_{\lambda}^{A_{0}}\left[\mathcal{P} f_{\lambda}^{A_{0}}-q\left(\mathcal{P} f_{\lambda}^{A_{0}}\right)^{2}-q\left(\pi w_{\lambda}^{A_{0}}\right)^{2}\right], \quad \delta f_{\lambda}^{A}=w_{\lambda}^{A}=g_{\lambda}^{A_{0}} w_{\lambda}^{A_{0}}  \tag{55}\\
\mathcal{P} \psi_{\lambda}^{A} & =g_{\lambda}^{A_{0}}\left[\mathcal{P} f_{\lambda}^{A_{0}} \phi_{\lambda}^{A_{0}}-q \mathcal{P} f_{\lambda}^{A_{0}} \xi_{\lambda}^{A_{0}}+\xi_{\lambda}^{A_{0}}-q\left(\mathcal{P} f_{\lambda}^{A_{0}}\right)^{2} \phi_{\lambda}^{A_{0}}-q\left(\pi w_{\lambda}^{A_{0}}\right)^{2} \phi_{\lambda}^{A_{0}}\right]  \tag{56}\\
\delta \psi_{\lambda}^{A} & =w_{\lambda}^{A} \phi_{\lambda}^{A}=g_{\lambda}^{A_{0}} w_{\lambda}^{A_{0}}\left(\phi_{\lambda}^{A_{0}}+q \xi_{\lambda}^{A_{0}}\right) \Rightarrow \phi_{\lambda}^{A}=\phi_{\lambda}^{A_{0}}+q \xi_{\lambda}^{A_{0}}  \tag{57}\\
\delta_{\lambda}^{A} & =w_{\lambda}^{A} \phi_{\lambda}^{A} \otimes \phi_{\lambda}^{A_{, *}}=g_{\lambda}^{A_{0}} w_{\lambda}^{A_{0}}\left(\phi_{\lambda}^{A_{0}}+q \xi_{\lambda}^{A}\right) \otimes\left(\phi_{\lambda}^{A_{0}, *}+q \xi_{\lambda}^{A_{0}, *}\right)  \tag{58}\\
g_{\lambda}^{A_{0}} & =\left[\left(1-q \mathcal{P} f_{\lambda}^{A_{0}}\right)^{2}+\left(q \pi w_{\lambda}^{A_{0}}\right)^{2}\right]^{-1} \tag{59}
\end{align*}
$$

If $\delta_{\lambda}^{A_{0}}$ is regular at the threshold of $\sigma\left(A_{0}\right)$ then an analysis of the threshold behavior of $\delta_{\lambda}^{A}$ is strongly controlled by the threshold behavior of $g_{\lambda}^{A_{0}}$.

Proof of Theorem (1) Part (3). One may straightforwardly take the presented techniques for rank-1 spectral shifts with the assignments $A_{0}=L_{0}, A=L$, $v_{\mathrm{N}}=v_{\mathrm{P}}=\chi_{0}$.

Of particular importance in our analysis will be the function

$$
\begin{equation*}
g_{\lambda}:=\left\{\left[1-q e^{-\lambda} \mathcal{P} E_{1}(-\lambda)\right]^{2}+\left[\pi q e^{-\lambda}\right]^{2}\right\}^{-1} \tag{60}
\end{equation*}
$$

which satisfies $w_{\lambda}^{L}=g_{\lambda} w_{\lambda}$. This function strongly controls the behavior of the spectral measure near the threshold due to the logarithmic divergence of $\mathcal{P} E_{1}(-\lambda)$ there.

## 4. Decay Estimates for $L_{0}$ and $L$

The Mourre estimate, see e.g. [23], has been proven for $L_{0}$ by Chen, Fröhlich, and Walcher [5], in order to prove its the spectrum is absolutely continuous and equal to $[0, \infty)$. We want to study pointwise decay estimates in time, which requires knowledge of the aysmptotic properties of the resolvent at thresholds. The Mourre estimates do not apply at thresholds so we will need to use alternative methods.

Local decay estimates for $L_{0}$ have been found by Durhuus and Gayral [9] in the context of more general noncommutative solitons (where $L_{0}$ corresponds their "diagonal case with 2 noncommuting spatial coordinates"). They found an unweighted estimate of the form $\left\|e^{-i t L_{0}} v\right\|_{\infty} \leq c|t|^{-1}(1+\log |t|)\|v\|_{1}$ for $|t| \geq 1$. We present an alternative approach, in the context of Jacobi operators, which enhances the local decay estimate for the free Schrödinger operator and provides integrable decay for rank one boundary perturbations thereof. We find weighted estimates

$$
\left\|W_{\kappa, \tau} e^{-i t L_{0}} W_{\kappa, \tau} v\right\|_{\infty}<c t^{-1}\|v\|_{1}, \quad\left\|W_{\kappa, \tau} e^{-i t L} \Pi_{\mathrm{e}}^{L} W_{\kappa, \tau} v\right\|_{\infty}=\mathcal{O}\left(t^{-1} \log ^{-2} t\right)
$$

for $t \nearrow \infty$.

### 4.1. Weighted estimates for spectral vectors.

Definition 4.1. Let $\mathbb{S}_{r}, \mathbb{D}_{r} \subset \mathbb{C}$ be respectively the circle and the disc of radius $r>0$ centered at the origin and $u \in \mathscr{T}$ a formal sequence for which there exist constants $r, c>0$ for which $|u(x)| r^{-x}<c$ for all $x \in \mathbb{Z}_{+}$. The generating function of $u$ is the function $\zeta(u, \cdot): \mathbb{S}_{r^{\prime}} \rightarrow \mathbb{C}$ defined by $\zeta(u, s):=\sum_{x=0}^{\infty} u(x) s^{x}$, where $r^{\prime}<r$. This permits the presentation of $u$ via

$$
\begin{equation*}
u(x)=\oint_{\gamma} \mathrm{d} s(2 \pi i s)^{-1} s^{-x} \zeta(u, s) \tag{61}
\end{equation*}
$$

where $\gamma$ is any positively oriented simple closed curve in $\mathbb{D}_{r}$ which encloses and does not pass through the origin.

The Laguerre polynomials have the well-known generating function [35]

$$
\begin{equation*}
\zeta\left(\phi_{\lambda}, s\right):=\sum_{x=0}^{\infty} \phi_{\lambda}(x) s^{x}=(1-s)^{-1} \exp \left[-(1-s)^{-1} s \lambda\right], \quad|s|<1 \tag{62}
\end{equation*}
$$

We will also employ the notion of a reduced generating function.

Definition 4.2. For a given generating function $\zeta(v, s)$ of a vector $v$ let the reduced generating function be the function $\widehat{\zeta}(v, s):=(1-s) \zeta(v, s)$.

For example we have that $\widehat{\zeta}\left(\phi_{\lambda}, s\right)=\exp \left[-(1-s)^{-1} s \lambda\right]$.

Definition 4.3. For $s \in \mathbb{C}$ we let

$$
\begin{align*}
& r:=|s|, \quad \theta:=\arg (s), \quad \widehat{s}:=(1-s)^{-1} s,  \tag{63}\\
& \epsilon:=1-r^{2}, \quad \widehat{\epsilon}:=-2^{-1}+16^{-1} \epsilon . \tag{64}
\end{align*}
$$

Should many variables $s_{j}$ be present we define $\epsilon_{j}:=1-\left|s_{j}\right|^{2}$ correspondingly for each.

Lemma 4.1. One has the representation

$$
\begin{equation*}
\xi_{z}(x)=\int_{0}^{\infty} \mathrm{d} \lambda e^{-\lambda}(\lambda-z)^{-1}\left[\phi_{\lambda}(x)-\phi_{z}(x)\right] \tag{65}
\end{equation*}
$$

where

$$
\begin{equation*}
K(\eta, z, s):=(\eta-z)^{-1}[\exp (-\widehat{s} \eta)-\exp (-\widehat{s} z)] \tag{66}
\end{equation*}
$$

for $\eta \in \mathbb{R}_{+}, z \in \mathbb{C}, s \in \mathbb{S}_{r}$.

Proof. By the spectral representation of $R_{z}^{L_{0}}$ it is the case that $\psi_{z}(x)=\int_{0}^{\infty} \mathrm{d} \lambda e^{-\lambda}(\lambda-z)^{-1} \phi_{\lambda}(x)$, where we have used the normalization condition $\phi_{\lambda}(0)=1, \forall \lambda \in \sigma\left(L_{0}\right)$.

We are primarily concerned with estimates of operators in generating function presentation. In such forms one finds line integrations over dummy complex variables,
$s_{j}$, with a priori separate sums for each $x_{j} \in \mathbb{Z}_{+}$. One is therefore permitted to make the associated contours dependent on $x_{j}$. We then will hereafter take

$$
\begin{array}{rlrl}
r & :=\left|s_{1}\right| \stackrel{\text { set }}{=} 1-(x+\kappa)^{-1}<1 \\
\Rightarrow \quad & \epsilon & =1-\left[1-(x+\kappa)^{-1}\right]^{2}=2(x+\kappa)^{-1}+\mathcal{O}\left((x+\kappa)^{-2}\right) \tag{68}
\end{array}
$$

where dependence on $x \in \mathbb{Z}_{+}$and $1<\kappa \in \mathbb{R}$ will always be suppressed. One may observe that for sufficiently large $\kappa$ one may make $\epsilon$ as small as one would like and therefore will be treated as a small parameter. One may also observe that

$$
\begin{equation*}
\left|(1-s)^{-1}\right| \leq x+\kappa, \quad|\hat{s}| \leq x+\kappa-1<x+\kappa \tag{69}
\end{equation*}
$$

as well as the crucial estimate

$$
\begin{equation*}
\left|s^{-x}\right|<e<3 \tag{70}
\end{equation*}
$$

One is then permitted to work with polynomially weighted spaces instead of exponentially weighted ones.

Lemma 4.2. One has that $|\exp (-\widehat{s} \lambda)| \leq \exp (-\widehat{\epsilon} \lambda)$ for sufficiently large $\kappa$.

Proof. We recall that we have universally chosen that $|s| \stackrel{\text { set }}{=} 1-(x+\kappa)^{-1}$. Let $l:=\Re \widehat{s}, m:=1-\cos \theta$, and $n:=r \cos \theta-r^{2}$. Let $\mathbb{H}_{\mathrm{L}}:=\{z \in \mathbb{C}: \Re z<0\}$ be the left half plane. With regard to $\widehat{s}$ we need only consider $l:=\Re \widehat{s}$. By inspection of $\widehat{s}$ one can see that for $s \in \mathbb{S}$ the supremum of $|\exp (-\widehat{s} \lambda)|$ should be found in $\mathbb{S}_{r} \cap \mathbb{H}_{\mathrm{L}}$. Furthermore, although for unrestricted $s$ one has $0 \leq m \leq 2$ it is the case that for restricted $s$ one has $1 \leq m \leq 2$ and thereby $m=\mathcal{O}(1)$. For $s \in \mathbb{S}_{r} \cap \mathbb{H}_{\mathrm{L}}$ it follows that

$$
\begin{align*}
r & =(1-\epsilon)^{1 / 2}=1-\frac{1}{2} \epsilon+\mathcal{O}\left(\epsilon^{2}\right)  \tag{71}\\
n & =r \cos \theta-r^{2}=-m-\frac{1}{2}(1-m) \epsilon+\mathcal{O}\left(\epsilon^{2}\right) \tag{72}
\end{align*}
$$

$$
\begin{align*}
l & =n(\epsilon-2 n)^{-1}  \tag{73}\\
& =-m-\frac{1}{2}(1-m) \epsilon+\mathcal{O}\left(\epsilon^{2}\right)\left[\epsilon-2-m-\frac{1}{2}(1-m) \epsilon+\mathcal{O}\left(\epsilon^{2}\right)\right]^{-1}  \tag{74}\\
& =-2^{-1}+(4 m)^{-1} \epsilon+\widehat{l}(\epsilon) \tag{75}
\end{align*}
$$

where $\widehat{l}(\epsilon)=\mathcal{O}\left(\epsilon^{2}\right)$. One then has

$$
\begin{align*}
|\exp (-\widehat{s} \lambda)| & \leq \sup _{\theta \in[\pi / 2,3 \pi / 2]} \exp [-l \lambda]  \tag{76}\\
& =\sup _{\theta \in[\pi / 2,3 \pi / 2]} \exp \left\{-\left[-2^{-1}+(4 m)^{-1} \epsilon+\widehat{l}(\epsilon)\right] \lambda\right\}  \tag{77}\\
& =\sup _{\theta \in[\pi / 2,3 \pi / 2]} \exp \left\{-\left[-2^{-1}+8^{-1} \epsilon+\widehat{l}(\epsilon)\right] \lambda\right\} . \tag{78}
\end{align*}
$$

$\kappa$ is a free parameter and as it increases $\epsilon$ will decrease monotonically. Therefore there must exist a large enough $\kappa$ so that $|\widehat{l}(\epsilon)| \leq(16)^{-1} \epsilon, \forall \theta$. One then has that

$$
\begin{align*}
\sup _{\theta_{1} \in[\pi / 2,3 \pi / 2]} \exp \left\{-\left[-2^{-1}+8^{-1} \epsilon+\widehat{l}(\epsilon)\right] \lambda\right\} & \leq \exp \left\{-\left[-2^{-1}+8^{-1} \epsilon-16^{-1} \epsilon\right] \lambda\right\}  \tag{79}\\
& =\exp \left[-\left(-2^{-1}+16^{-1} \epsilon\right) \lambda\right] \tag{80}
\end{align*}
$$

Lemma 4.3. One has that $\left(\epsilon_{1}+\epsilon_{2}\right)^{-1}<4^{-1}\left(x_{1}+\kappa\right)\left(x_{2}+\kappa\right)$.

## Proof.

$$
\begin{align*}
\epsilon & =1-r^{2}=\left[2-(x+\kappa)^{-1}\right](x+\kappa)^{-1}  \tag{81}\\
\left(\epsilon_{1}+\epsilon_{2}\right)^{-1} & =\left(x_{1}+\kappa\right)\left(x_{2}+\kappa\right)\left[4-\left(x_{1}+\kappa\right)^{-1}-\left(x_{2}+\kappa\right)^{-1}\right]  \tag{82}\\
& <4^{-1}\left(x_{1}+\kappa\right)\left(x_{2}+\kappa\right) . \tag{83}
\end{align*}
$$

Lemma 4.4. One has the generating function representation

$$
\begin{equation*}
\xi_{z}(x)=\oint_{\mathbb{S}_{r}} \mathrm{~d} s(2 \pi i s)^{-1} s^{-x} \zeta\left(\xi_{z}, s\right), \quad \forall z \in \mathbb{C} \tag{84}
\end{equation*}
$$

where $\zeta\left(\xi_{z}, s\right):=(1-s)^{-1} \int_{0}^{\infty} \mathrm{d} \eta e^{-\eta} K(\eta, z, s)$.

Proof. First, let $z \in \mathbb{C} \backslash \mathbb{R}_{+}=: \Sigma$ and $l:=\Re \widehat{s}$. Since $|s|<1$ there exists a $c>0$ such that $|\widehat{s}|<c$. It follows that

$$
\begin{align*}
|K(\eta, z, s)| & \leq\left|(\eta-z)^{-1}\right|[|\exp (-\widehat{s} \eta)|+|\exp (-\widehat{s} z)|]  \tag{85}\\
& \leq[\operatorname{dist}(\Sigma, z)]^{-1}[\exp (-l \eta)+\exp (-c|z|)]<\infty \tag{86}
\end{align*}
$$

Second, let $z \equiv \lambda \in \mathbb{R}_{+}$. By mean value theorem one has

$$
\begin{align*}
K(\eta, \lambda, s)= & (\eta-\lambda)^{-1}[\Re \exp (-\widehat{s} \eta)-\Re \exp (-\widehat{s} z)]  \tag{87}\\
& +i(\eta-\lambda)^{-1}[\Im \exp (-\widehat{s} \eta)-\Im \exp (-\widehat{s} z)]  \tag{88}\\
= & \mathrm{d}_{\eta} \Re \exp (-\widehat{s} \eta) L_{\eta=\mu_{1}}+i \mathrm{~d}_{\eta} \Im \exp (-\widehat{s} \eta)\left\lfloor_{\eta=\mu_{2}}\right.  \tag{89}\\
= & \frac{1}{2}\left[(-\widehat{s}) \exp \left(-\widehat{s} \mu_{1}\right)+(-\overline{\widehat{s}}) \exp \left(-\overline{\widehat{s}} \mu_{1}\right)\right.  \tag{90}\\
& \left.+(-\widehat{s}) \exp \left(-\widehat{s} \mu_{2}\right)-(-\overline{\widehat{s}}) \exp \left(-\overline{\widehat{s}} \mu_{2}\right)\right] \tag{91}
\end{align*}
$$

where $\mu_{j} \equiv \mu_{j}(r, \theta, \eta, \lambda) \in[\min (\eta, \lambda), \max (\eta, \lambda)]$. Then

$$
\begin{align*}
|K(\eta, \lambda, s)| \leq & \frac{1}{2}\left[|\widehat{s}|\left|\exp \left(-\widehat{s} \mu_{1}\right)\right|+|\widehat{s}|\left|\exp \left(-\widehat{s} \mu_{1}\right)\right|\right.  \tag{92}\\
& \left.+|\widehat{s}|\left|\exp \left(-\widehat{s} \mu_{2}\right)\right|+|\widehat{s}|\left|\exp \left(-\widehat{s} \mu_{2}\right)\right|\right]  \tag{93}\\
= & |\widehat{s}|\left[\exp \left(-l \mu_{1}\right)+\exp \left(-l \mu_{2}\right)\right]  \tag{94}\\
\leq & 2|\widehat{s}| \exp [-\widehat{\epsilon}(\eta+\lambda)]<\infty \tag{95}
\end{align*}
$$

One may observe that

$$
\begin{align*}
\int_{0}^{\infty} \mathrm{d} \eta e^{-\eta}|K(\eta, z, s)| & \leq \int_{0}^{\infty} \mathrm{d} \eta e^{-\eta} 2|\widehat{s}| \exp [-\widehat{\epsilon}(\eta+\lambda)]  \tag{96}\\
& =2|\widehat{s}| \exp (-\widehat{\epsilon} \lambda) \int_{0}^{\infty} \mathrm{d} \eta \exp [-(1+\widehat{\epsilon}) \eta]  \tag{97}\\
& =2|\widehat{s}| \exp (-\widehat{\epsilon} \lambda)(1+\widehat{\epsilon})^{-1}<\infty \tag{98}
\end{align*}
$$

The multi-integral of the generating function representation of $\xi_{z}(x)$ converges absolutely and thereby Fubini's theorom permits

$$
\begin{align*}
\xi_{z}(x) & =\int_{0}^{\infty} \mathrm{d} \eta e^{-\eta}(\eta-z)^{-1}\left[\phi_{\eta}(x)-\phi_{z}(x)\right]  \tag{99}\\
& =\int_{0}^{\infty} \mathrm{d} \eta e^{-\eta} \oint_{\mathbb{S}_{r}} \mathrm{~d} s(2 \pi i s)^{-1} s^{-x}(1-s)^{-1} K(\eta, z, s)  \tag{100}\\
& =\oint_{\mathbb{S}_{r}} \mathrm{~d} s(2 \pi i s)^{-1} s^{-x}(1-s)^{-1} \int_{0}^{\infty} \mathrm{d} \eta e^{-\eta} K(\eta, z, s) \tag{101}
\end{align*}
$$

for all $z \in \mathbb{C}$.

Lemma 4.5. It is the case that

$$
\begin{equation*}
\left|\mathrm{d}_{\lambda}^{n} \widehat{\zeta}\left(\phi_{\lambda}, s\right)\right|<\widehat{c}\left(\phi_{\lambda}, n\right) \exp (-\widehat{\epsilon} \lambda), \quad n \in \mathbb{Z}_{+} \tag{102}
\end{equation*}
$$

where $\widehat{c}\left(\phi_{\lambda}, n\right):=(x+\kappa)^{n}$ and

$$
\begin{equation*}
\left|\mathrm{d}_{\lambda}^{n} \widehat{\zeta}\left(\xi_{\lambda}, s\right)\right|<\widehat{c}\left(\xi_{\lambda}, n\right) \exp (-\widehat{\epsilon} \lambda), \quad n=0,1,2 \tag{103}
\end{equation*}
$$

where

$$
\begin{equation*}
\widehat{c}\left(\xi_{\lambda}, 0\right):=4(x+\kappa), \quad \widehat{c}\left(\xi_{\lambda}, 1\right):=6(x+\kappa), \quad \widehat{c}\left(\xi_{\lambda}, 2\right):=8(x+\kappa)^{2} . \tag{104}
\end{equation*}
$$

Proof. For $\phi_{\lambda}$ :

$$
\begin{align*}
\left|\mathrm{d}_{\lambda}^{n} \widehat{\zeta}\left(\phi_{\lambda}, s\right)\right| & =\left|\mathrm{d}_{\lambda}^{n} \exp (-\widehat{s} \lambda)\right|=\left|\widehat{s}^{n} \exp (-\widehat{s} \lambda)\right|  \tag{105}\\
& \leq|\widehat{s}|^{n}|\exp (-\widehat{s} \lambda)| \leq|\widehat{s}|^{n} \exp (-\widehat{\epsilon} \lambda)<(x+\kappa)^{n} \exp (-\widehat{\epsilon} \lambda) \tag{106}
\end{align*}
$$

For $\xi_{\lambda}$ : One may observe that $K(\eta, \lambda, s)=K(\lambda, \eta, s)$. Then, by integration by parts, one has

$$
\begin{align*}
\mathrm{d}_{\lambda} \widehat{\zeta}\left(\xi_{\lambda}, s\right) & =\int_{0}^{\infty} \mathrm{d} \eta e^{-\eta} \mathrm{d}_{\lambda} K(\eta, \lambda, s)=\int_{0}^{\infty} \mathrm{d} \eta e^{-\eta} \mathrm{d}_{\eta} K(\eta, \lambda, s)  \tag{107}\\
& =-K(0, \lambda, s)+\widehat{\zeta}\left(\xi_{\lambda}, s\right) \tag{108}
\end{align*}
$$

and thereby

$$
\begin{equation*}
\mathrm{d}_{\lambda}^{n} \widehat{\zeta}\left(\xi_{\lambda}, s\right)=-\sum_{k=0}^{n-1} \mathrm{~d}_{\lambda} K(0, \lambda, s)+\widehat{\zeta}\left(\xi_{\lambda}, s\right) \tag{109}
\end{equation*}
$$

where the sum is defined to vanish when the upper bound is negative. It follows that

$$
\begin{equation*}
|K(\eta, \lambda, s)| \leq 2|\widehat{s}| \exp [-\widehat{\epsilon}(\eta+\lambda)], \quad\left|\widehat{\zeta}\left(\xi_{\lambda}, s\right)\right| \leq 2|\widehat{s}|(1+\widehat{\epsilon})^{-1} \exp (-\widehat{\epsilon} \lambda) \tag{110}
\end{equation*}
$$

Consider an arbitrary $f \in C^{2}(\mathbb{R}, \mathbb{R})$ and let $f_{*}$ be its Newton quotient so that

$$
\begin{equation*}
f_{*}\left(a_{0}, a\right):=\left(a_{0}-a\right)^{-1}\left[f\left(a_{0}\right)-f(a)\right] . \tag{111}
\end{equation*}
$$

One has by mean value theorem

$$
\begin{align*}
& \begin{aligned}
& \mathrm{d}_{a} f_{*}\left(a_{0}, a\right)=\left(a_{0}-a\right)^{-2}\left[f\left(a_{0}\right)-f(a)-\left(a_{0}-a\right) \mathrm{d}_{a} f(a)\right] \\
&=\left(a_{0}-a\right)^{-1}\left[\mathrm{~d}_{a_{1}} f\left(a_{1}\right)-\mathrm{d}_{a} f(a)\right], \quad a_{1} \in\left[\min \left(a_{0}, a\right), \max \left(a_{0}, a\right)\right] \\
&=\left(a_{0}-a\right)^{-1}\left(a_{1}-a\right) \mathrm{d}_{a_{2}}^{2} f\left(a_{2}\right), \quad a_{2} \in\left[\min \left(a_{1}, a\right), \max \left(a_{1}, a\right)\right] \\
& \Rightarrow \quad\left|\mathrm{d}_{a} f_{*}\left(a_{0}, a\right)\right| \leq\left|\left(a_{0}-a\right)^{-1}\right|\left|a_{1}-a\right|\left|\mathrm{d}_{a_{2}}^{2} f\left(a_{2}\right)\right| \leq\left|\mathrm{d}_{a_{2}}^{2} f\left(a_{2}\right)\right|
\end{aligned} . \tag{112}
\end{align*}
$$

Let $(\Re, \Im) z$ be a presentation for the real and imaginary parts of $z \in \mathbb{C}$ whose ordering in compatible with the respective ordering of $\pm$. Let $i_{+}:=1, i_{-}:=i$ and
$\mu_{ \pm} \in[\min (\eta, \lambda), \max (\eta, \lambda)]$. It follows that

$$
\begin{align*}
\left|\mathrm{d}_{\lambda}(\Re, \Im) K(\eta, \lambda, s)\right| & \leq\left.\left|\mathrm{d}_{\lambda}^{2}(\Re, \Im) \exp (-\widehat{s} \lambda)\right|\right|_{\lambda=\mu_{ \pm}}  \tag{116}\\
& =\left.\left|\mathrm{d}_{\lambda}^{2}\left(2 i_{ \pm}\right)^{-1}[\exp (-\widehat{s} \lambda) \pm \exp (-\widehat{\widehat{s}} \lambda)]\right|\right|_{\lambda=\mu_{ \pm}}  \tag{117}\\
& \leq|\widehat{s}|^{2} \exp [-\widehat{\epsilon}(\eta+\lambda)]  \tag{118}\\
\Rightarrow \quad\left|\mathrm{d}_{\lambda} K(\eta, \lambda, s)\right| & \leq 2|\widehat{s}|^{2} \exp [-\widehat{\epsilon}(\eta+\lambda)] \tag{119}
\end{align*}
$$

Then

$$
\begin{align*}
\left|\widehat{\zeta}\left(\xi_{\lambda}, s\right)\right| & \leq 2|\widehat{s}|(1+\widehat{\epsilon})^{-1} \exp (-\widehat{\epsilon} \lambda)<4(x+\kappa) \exp (-\widehat{\epsilon} \lambda)  \tag{120}\\
\left|\mathrm{d}_{\lambda} \widehat{\zeta}\left(\xi_{\lambda}, s\right)\right| & \leq 2|\widehat{s}|\left[(1+\widehat{\epsilon})^{-1}+1\right] \exp (-\widehat{\epsilon} \lambda)<6(x+\kappa) \exp (-\widehat{\epsilon} \lambda)  \tag{121}\\
\left|\mathrm{d}_{\lambda} \widehat{\zeta}\left(\xi_{\lambda}, s\right)\right| & \left.\leq 2|\widehat{s}|\left[(1+\widehat{\epsilon})^{-1}+1+\mid \widehat{s}\right]\right] \exp (-\widehat{\epsilon} \lambda)<8(x+\kappa)^{2} \exp (-\widehat{\epsilon} \lambda) \tag{122}
\end{align*}
$$

REMARK. If estimates of $\mathrm{d}_{\lambda}^{n} \widehat{\zeta}\left(\xi_{\lambda}, s\right)$ for $2<n \in \mathbb{Z}$ were required the above method would not follow so straightforwardly due to the inapplicability of the mean value theorem for yet higher derivatives.

Corollary 4.1. One has that

$$
\begin{equation*}
\left|\mathrm{d}_{\lambda}^{n}\left[w_{\lambda}^{1 / 2} \phi_{\lambda}(x)\right]\right|<c\left(\phi_{\lambda}, n\right) \exp \left(-16^{-1} \epsilon \lambda\right), \quad n=0,1,2 \tag{123}
\end{equation*}
$$

where $c\left(\phi_{\lambda}, n\right):=3^{n+1}(x+\kappa)^{n+1}$ and

$$
\begin{equation*}
\left|\mathrm{d}_{\lambda}^{n}\left[w_{\lambda}^{1 / 2} \xi_{\lambda}(x)\right]\right|<c\left(\xi_{\lambda}, n\right) \exp \left(-16^{-1} \epsilon \lambda\right), \quad n=0,1,2 \tag{124}
\end{equation*}
$$

where

$$
\begin{equation*}
c\left(\xi_{\lambda}, 0\right):=12(x+\kappa)^{2}, \quad c\left(\xi_{\lambda}, 1\right):=24(x+\kappa)^{2}, \quad c\left(\xi_{\lambda}, 2\right):=36(x+\kappa)^{3} \tag{125}
\end{equation*}
$$

Proof. For $\phi_{\lambda}$ :

$$
\begin{align*}
&\left|\mathrm{d}_{\lambda}^{n}\left[w_{\lambda}^{1 / 2} \phi_{\lambda}(x)\right]\right|=\left|\oint_{\mathbb{S}_{r}} \mathrm{~d} s(2 \pi i s)^{-1} s^{-x}(1-s)^{-1} \mathrm{~d}_{\lambda}^{n}\left[w_{\lambda}^{1 / 2} \widehat{\zeta}\left(\phi_{\lambda}, s\right)\right]\right|  \tag{126}\\
&= \mid \oint_{\mathbb{S}_{r}} \mathrm{~d} s(2 \pi i s)^{-1} s^{-x}(1-s)^{-1}  \tag{127}\\
& \left.\times \sum_{k=0}^{n}\binom{n}{k} \mathrm{~d}_{\lambda}^{n-k} w_{\lambda}^{1 / 2} \mathrm{~d}_{\lambda}^{k} \widehat{\zeta}\left(\phi_{\lambda}, s\right) \right\rvert\,  \tag{128}\\
& \leq \oint_{\mathbb{S}_{r}}\left|\mathrm{~d} s(2 \pi i s)^{-1}\right|\left|s^{-x}\right|\left|(1-s)^{-1}\right|  \tag{129}\\
& \times \sum_{k=0}^{n}\binom{n}{k}\left|\mathrm{~d}_{\lambda}^{n-k} w_{\lambda}^{1 / 2}\right|\left|\mathrm{d}_{\lambda}^{k} \widehat{\zeta}\left(\phi_{\lambda}, s\right)\right|  \tag{130}\\
&<(1)(3)(x+\kappa) \sum_{k=0}^{n}\binom{n}{k} 2^{-(n-k)} w_{\lambda}^{1 / 2} \widehat{c}\left(\phi_{\lambda}, k\right) \exp (-\widehat{\epsilon} \lambda)  \tag{131}\\
&= 3(x+\kappa) \sum_{k=0}^{n}\binom{n}{k} 2^{-(n-k)} \widehat{c}\left(\phi_{\lambda}, k\right) \exp \left(-16^{-1} \epsilon \lambda\right) \tag{132}
\end{align*}
$$

For $n=0$ :

$$
\begin{equation*}
\left|w_{\lambda}^{1 / 2} \phi_{\lambda}(x)\right|<3(x+\kappa) \widehat{c}\left(\phi_{\lambda}, 0\right) \exp \left(-16^{-1} \epsilon \lambda\right)=3(x+\kappa) \exp \left(-16^{-1} \epsilon \lambda\right) \tag{133}
\end{equation*}
$$

For $n=1$ :

$$
\begin{align*}
\left|\mathrm{d}_{\lambda}\left[w_{\lambda}^{1 / 2} \phi_{\lambda}(x)\right]\right| & <3(x+\kappa)\left[2^{-1} \widehat{c}\left(\phi_{\lambda}, 0\right)+\widehat{c}\left(\phi_{\lambda}, 1\right)\right] \exp \left(-16^{-1} \epsilon \lambda\right)  \tag{134}\\
& =3(x+\kappa)\left[2^{-1}+(x+\kappa)\right] \exp \left(-16^{-1} \epsilon \lambda\right)  \tag{135}\\
& <6(x+\kappa)^{2} \exp \left(-16^{-1} \epsilon \lambda\right) . \tag{136}
\end{align*}
$$

For $n=2$ :

$$
\begin{align*}
&\left|\mathrm{d}_{\lambda}^{2}\left[w_{\lambda}^{1 / 2} \phi_{\lambda}(x)\right]\right|< 3(x+\kappa)\left[2^{-2} \widehat{c}\left(\phi_{\lambda}, 0\right)+2^{-1} \widehat{c}\left(\phi_{\lambda}, 1\right)+\widehat{c}\left(\phi_{\lambda}, 2\right)\right]  \tag{137}\\
& \quad \times \exp \left(-16^{-1} \epsilon \lambda\right)  \tag{138}\\
&= 3(x+\kappa)\left[2^{-2}+2^{-1}(x+\kappa)+(x+\kappa)^{2}\right]  \tag{139}\\
& \quad \times \exp \left(-16^{-1} \epsilon \lambda\right)  \tag{140}\\
&< 9(x+\kappa)^{3} \exp \left(-16^{-1} \epsilon \lambda\right) \tag{141}
\end{align*}
$$

For $\xi_{\lambda}$ :

$$
\begin{equation*}
\left|w_{\lambda}^{1 / 2} \xi_{\lambda}(x)\right|<3(x+\kappa) \sum_{k=0}^{n}\binom{n}{k} 2^{-(n-k)} \widehat{c}\left(\xi_{\lambda}, k\right) \exp \left(-16^{-1} \epsilon \lambda\right) \tag{142}
\end{equation*}
$$

For $n=0$ :

$$
\begin{equation*}
\left|w_{\lambda}^{1 / 2} \xi_{\lambda}(x)\right|<3(x+\kappa) \widehat{c}\left(\xi_{\lambda}, 0\right) \exp \left(-16^{-1} \epsilon \lambda\right)=12(x+\kappa)^{2} \exp \left(-16^{-1} \epsilon \lambda\right) \tag{143}
\end{equation*}
$$

For $n=1$ :

$$
\begin{align*}
\left|\mathrm{d}_{\lambda}\left[w_{\lambda}^{1 / 2} \xi_{\lambda}(x)\right]\right| & <3(x+\kappa)\left[2^{-1} \widehat{c}\left(\xi_{\lambda}, 0\right)+\widehat{c}\left(\xi_{\lambda}, 1\right)\right] \exp \left(-16^{-1} \epsilon \lambda\right)  \tag{144}\\
& =3(x+\kappa)[2(x+\kappa)+6(x+\kappa)] \exp \left(-16^{-1} \epsilon \lambda\right)  \tag{145}\\
& =24(x+\kappa)^{2} \exp \left(-16^{-1} \epsilon \lambda\right) \tag{146}
\end{align*}
$$

For $n=2$ :

$$
\begin{align*}
\left|\mathrm{d}_{\lambda}^{2}\left[w_{\lambda}^{1 / 2} \xi_{\lambda}(x)\right]\right|< & 3(x+\kappa)\left[2^{-2} \widehat{c}\left(\xi_{\lambda}, 0\right)+2^{-1} \widehat{c}\left(\xi_{\lambda}, 1\right)+\widehat{c}\left(\xi_{\lambda}, 2\right)\right]  \tag{147}\\
& \times \exp \left(-16^{-1} \epsilon \lambda\right)  \tag{148}\\
= & 3(x+\kappa)\left[(x+\kappa)+3(x+\kappa)+8(x+\kappa)^{2}\right]  \tag{149}\\
& \times \exp \left(-16^{-1} \epsilon \lambda\right)  \tag{150}\\
< & 36(x+\kappa)^{3} \exp \left(-16^{-1} \epsilon \lambda\right) \tag{151}
\end{align*}
$$

### 4.2. Local time decay for $L_{0}$.

Proof of Theorem (2). Let $t>0,\left|s_{j}\right| \equiv r_{j} \stackrel{\text { set }}{=} 1-\left(x_{j}+\kappa\right)^{-1}, j=1,2$, for $1<\kappa \in \mathbb{R}$. It is the case that

$$
\begin{align*}
& \left|e^{-i t L_{0}}\left(x_{1}, x_{2}\right)\right|=\left|\int_{0}^{\infty} \mathrm{d} \lambda e^{-i t \lambda} w_{\lambda} \phi_{\lambda}\left(x_{1}\right) \phi_{\lambda}\left(x_{2}\right)\right|  \tag{152}\\
& =\left|\int_{0}^{\infty} \mathrm{d} \lambda(-i t)^{-1} \mathrm{~d}_{\lambda} e^{-i t \lambda}\left[w_{\lambda}^{1 / 2} \phi_{\lambda}\left(x_{1}\right)\right]\left[w_{\lambda}^{1 / 2} \phi_{\lambda}\left(x_{2}\right)\right]\right|  \tag{153}\\
& \leq \mid-(-i t)^{-1}-\int_{0}^{\infty} \mathrm{d} \lambda(-i t)^{-1} e^{-i t \lambda}\left\{\mathrm{~d}_{\lambda}\left[w_{\lambda}^{1 / 2} \phi_{\lambda}\left(x_{1}\right)\right]\right.  \tag{154}\\
& \left.\times\left[w_{\lambda}^{1 / 2} \phi_{\lambda}\left(x_{2}\right)\right]+\left[w_{\lambda}^{1 / 2} \phi_{\lambda}\left(x_{1}\right)\right] \mathrm{d}_{\lambda}\left[w_{\lambda}^{1 / 2} \phi_{\lambda}\left(x_{2}\right)\right]\right\} \mid  \tag{155}\\
& \leq t^{-1}\left(1+\int_{0}^{\infty} \mathrm{d} \lambda\left\{\left|\mathrm{~d}_{\lambda}\left[w_{\lambda}^{1 / 2} \phi_{\lambda}\left(x_{1}\right)\right]\right|\left|w_{\lambda}^{1 / 2} \phi_{\lambda}\left(x_{2}\right)\right|\right.\right.  \tag{156}\\
& \left.\left.+\left|w_{\lambda}^{1 / 2} \phi_{\lambda}\left(x_{1}\right)\right|\left|\mathrm{d}_{\lambda}\left[w_{\lambda}^{1 / 2} \phi_{\lambda}\left(x_{2}\right)\right]\right|\right\}\right)  \tag{157}\\
& <t^{-1}\left(1+\int_{0}^{\infty} \mathrm{d} \lambda\left\{(6)\left(x_{1}+\kappa\right)^{2} \exp \left(-16^{-1} \epsilon_{1} \lambda\right)(3)\left(x_{2}+\kappa\right) \exp \left(-16^{-1} \epsilon_{2} \lambda\right)\right.\right.  \tag{158}\\
& \left.\left.+(3)\left(x_{1}+\kappa\right) \exp \left(-16^{-1} \epsilon_{1} \lambda\right)(6)\left(x_{2}+\kappa\right) \exp \left(-16^{-1} \epsilon_{2} \lambda\right)\right\}\right)  \tag{159}\\
& \leq t^{-1}\left\{1+18\left(x_{1}+\kappa\right)^{2}\left(x_{2}+\kappa\right)^{2} \int_{0}^{\infty} \mathrm{d} \lambda \exp \left[-16^{-1}\left(\epsilon_{1}+\epsilon_{2}\right) \lambda\right]\right\}  \tag{160}\\
& =t^{-1}\left[1+288\left(x_{1}+\kappa\right)^{2}\left(x_{2}+\kappa\right)^{2}\left(\epsilon_{1}+\epsilon_{2}\right)^{-1}\right]  \tag{161}\\
& <t^{-1}\left[1+72\left(x_{1}+\kappa\right)^{3}\left(x_{2}+\kappa\right)^{3}\right]  \tag{162}\\
& \leq 73\left(x_{1}+\kappa\right)^{3}\left(x_{2}+\kappa\right)^{3} t^{-1} . \tag{163}
\end{align*}
$$

4.3. Local time decay for $L$. We recall without proof Lemma 3.12 from [21]:

Lemma. Let $\mathscr{B}$ be a Banach space and $\lambda_{+}>\lambda_{-}$be real constants. If $F(\lambda)$ has the properties
(1) $F \in C\left(\lambda_{-}, \lambda_{+} ; \mathscr{B}\right)$
(2) $F\left(\lambda_{-}\right)=F(\lambda)=0, \quad \lambda>\lambda_{+}$
(3) $\mathrm{d}_{\lambda} F \in L^{1}\left(\lambda_{-}+\delta, \lambda_{+} ; \mathscr{B}\right), \quad \forall \delta>0$
(4) $\mathrm{d}_{\lambda} F(\lambda)=\mathcal{O}\left(\left[\lambda-\lambda_{-}\right]^{-1} \log ^{-3}\left[\lambda-\lambda_{-}\right]\right), \quad \lambda \searrow \lambda_{-}$
(5) $\mathrm{d}_{\lambda}^{2} F(\lambda)=\mathcal{O}\left(\left[\lambda-\lambda_{-}\right]^{-2} \log ^{-2}\left[\lambda-\lambda_{-}\right]\right), \quad \lambda \searrow \lambda_{-}$
then

$$
\begin{equation*}
\int_{\lambda_{-}}^{\infty} \mathrm{d} \lambda e^{-i t \lambda} F(\lambda)=\mathcal{O}\left(t^{-1} \log ^{-2} t\right), \quad t \nearrow \infty \tag{164}
\end{equation*}
$$

in the norm of $\mathscr{B}$.

Proof of Theorem (3). Let $\mathscr{B}=\left\{A \in \mathcal{L}(\mathscr{T}):\|A\|_{\mathscr{B}}<\infty\right\}$ be the Banach space complete in the norm

$$
\begin{equation*}
\|A\|_{\mathscr{B}}:=\sup _{v \in \ell^{1}} \frac{\left\|W_{\kappa, \tau} A W_{\kappa, \tau} v\right\|_{1}}{\|v\|_{1}} . \tag{165}
\end{equation*}
$$

Let $F(\lambda)=\delta_{\lambda}^{L}$. We will verify the appropriate properties of $F(\lambda)$ for $\lambda_{-}=0$ and $\lambda_{+}=\infty$.

We recall that

$$
\begin{equation*}
F\left(\lambda, x_{1}, x_{2}\right)=w_{\lambda}^{L} \phi_{\lambda}^{L}\left(x_{1}\right) \phi_{\lambda}^{L}\left(x_{2}\right)=g_{\lambda} w_{\lambda}\left[\phi_{\lambda}\left(x_{1}\right)+q \xi_{\lambda}\left(x_{1}\right)\right]\left[\phi_{\lambda}\left(x_{2}+q \xi_{\lambda}\left(x_{2}\right)\right]\right. \tag{166}
\end{equation*}
$$

One may observe that

$$
\begin{align*}
\left|\mathrm{d}_{\lambda}^{n}\left[w_{\lambda}^{1 / 2} \phi_{\lambda}^{L}(x)\right]\right| & \leq\left|\mathrm{d}_{\lambda}^{n}\left[w_{\lambda}^{1 / 2} \phi_{\lambda}(x)\right]\right|+q\left|\mathrm{~d}_{\lambda}^{n}\left[w_{\lambda}^{1 / 2} \xi_{\lambda}(x)\right]\right|  \tag{167}\\
& <\left[c\left(\phi_{\lambda}, n\right)+q c\left(\xi_{\lambda}, n\right)\right] \exp \left(-16^{-1} \epsilon \lambda\right)  \tag{168}\\
& <c\left(\phi_{\lambda}^{L}, n\right) \exp \left(-16^{-1} \epsilon \lambda\right), \quad n=0,1,2 \tag{169}
\end{align*}
$$

where here we choose

$$
\begin{align*}
c\left(\phi_{\lambda}^{L}, 0\right) & :=3(1+4 q)(x+\kappa)^{2},  \tag{170}\\
c\left(\phi_{\lambda}^{L}, 1\right) & :=6(1+4 q)(x+\kappa)^{2},  \tag{171}\\
c\left(\phi_{\lambda}^{L}, 2\right) & :=9(1+4 q)(x+\kappa)^{3} . \tag{172}
\end{align*}
$$

One may see by inspection that $g_{\lambda}:=\left\{\left[1-q e^{-\lambda} \mathcal{P} E_{1}(-\lambda)\right]^{2}+\left[\pi q e^{-\lambda}\right]^{2}\right\}^{-1}$ has the properties:

$$
\begin{align*}
g_{\lambda} & =\left|g_{\lambda}\right| \leq \widehat{g}_{0}(q)<\infty, \quad \forall \lambda \in[0, \infty)  \tag{173}\\
\left|\mathrm{d}_{\lambda} g_{\lambda}\right| & \leq \widehat{g}_{0}(q) \widehat{g}_{1}(q, \delta)<\infty, \quad \forall \lambda \in[\delta, \infty)  \tag{174}\\
g_{0} & =g_{\infty}=0,  \tag{175}\\
\mathrm{~d}_{\lambda} g_{\lambda} & =\mathcal{O}\left(\lambda^{-1} \log ^{-1} \lambda\right), \quad \lambda \searrow 0  \tag{176}\\
\mathrm{~d}_{\lambda}^{2} g_{\lambda} & =\mathcal{O}\left(\lambda^{-2} \log ^{-3} \lambda\right), \quad \lambda \searrow 0  \tag{177}\\
& \subset \mathcal{O}\left(\lambda^{-2} \log ^{-2} \lambda\right) \tag{178}
\end{align*}
$$

where $0<\widehat{g}_{0}(q), \widehat{g}_{1}(q, \delta)<\infty$ are constants whose other properties are not needed here. $g_{\lambda}$ is the only function of $\lambda$ involved in the definition of $F(\lambda)$ whose derivatives are unbounded in the neighborhood of the threshold $\lambda=0$ and thereby the derivatives of $g_{\lambda}$ are dominant in determining the properties of the derivatives of $F(\lambda)$.

Properties (1), (2): One may observe that the properties follow by inspection.

Property (3): For $\lambda \in[\delta, \infty)$ one has

$$
\begin{align*}
\left|\mathrm{d}_{\lambda} F\left(\lambda, x_{1}, x_{2}\right)\right|= & \left|\mathrm{d}_{\lambda}\left\{g_{\lambda}\left[w_{\lambda}^{1 / 2} \phi_{\lambda}^{L}\left(x_{1}\right)\right]\left[w_{\lambda}^{1 / 2} \phi_{\lambda}^{L}\left(x_{2}\right)\right]\right\}\right|  \tag{179}\\
\leq & \left|\mathrm{d}_{\lambda} g_{\lambda}\right|\left|\left[w_{\lambda}^{1 / 2} \phi_{\lambda}^{L}\left(x_{1}\right)\right]\right|\left|\left[w_{\lambda}^{1 / 2} \phi_{\lambda}^{L}\left(x_{2}\right)\right]\right|  \tag{180}\\
& +\left|g_{\lambda}\right|\left|\mathrm{d}_{\lambda}\left[w_{\lambda}^{1 / 2} \phi_{\lambda}^{L}\left(x_{1}\right)\right]\right|\left|\left[w_{\lambda}^{1 / 2} \phi_{\lambda}^{L}\left(x_{2}\right)\right]\right|  \tag{181}\\
& +\left|g_{\lambda}\right|\left|\left[w_{\lambda}^{1 / 2} \phi_{\lambda}^{L}\left(x_{1}\right)\right]\right|\left|\mathrm{d}_{\lambda}\left[w_{\lambda}^{1 / 2} \phi_{\lambda}^{L}\left(x_{2}\right)\right]\right|  \tag{182}\\
< & \widehat{g}_{0}(q) \widehat{g}_{1}(q, \delta)(3)(1+4 q)\left(x_{1}+\kappa\right)^{2} \exp \left(-16^{-1} \epsilon_{1} \lambda\right)  \tag{183}\\
& \times(3)(1+4 q)\left(x_{2}+\kappa\right)^{2} \exp \left(-16^{-1} \epsilon_{1} \lambda\right)  \tag{184}\\
& +\widehat{g}_{0}(q)(6)(1+4 q)\left(x_{1}+\kappa\right)^{2} \exp \left(-16^{-1} \epsilon_{1} \lambda\right)  \tag{185}\\
& \times(3)(1+4 q)\left(x_{2}+\kappa\right)^{2} \exp \left(-16^{-1} \epsilon_{2} \lambda\right)  \tag{186}\\
& +\widehat{g}_{0}(q)(3)(1+4 q)\left(x_{1}+\kappa\right)^{2} \exp \left(-16^{-1} \epsilon_{1} \lambda\right)  \tag{187}\\
& \times(6)(1+4 q)\left(x_{2}+\kappa\right)^{2} \exp \left(-16^{-1} \epsilon_{2} \lambda\right)  \tag{188}\\
= & c_{0}(q, \delta)\left(x_{1}+\kappa\right)^{2}\left(x_{2}+\kappa\right)^{2} \exp \left[-16^{-1}\left(\epsilon_{1}+\epsilon_{2}\right) \lambda\right] \tag{189}
\end{align*}
$$

where $c_{0}(q, \delta)$ is a constant.

Property (4): For $\lambda \searrow 0$ one has

$$
\begin{align*}
\left|\mathrm{d}_{\lambda} F\left(\lambda, x_{1}, x_{2}\right)\right|= & \left|\mathrm{d}_{\lambda}\left\{g_{\lambda}\left[w_{\lambda}^{1 / 2} \phi_{\lambda}^{L}\left(x_{1}\right)\right]\left[w_{\lambda}^{1 / 2} \phi_{\lambda}^{L}\left(x_{2}\right)\right]\right\}\right|  \tag{190}\\
\leq & \left|\mathrm{d}_{\lambda} g_{\lambda}\right|\left|\left[w_{\lambda}^{1 / 2} \phi_{\lambda}^{L}\left(x_{1}\right)\right]\right|\left|\left[w_{\lambda}^{1 / 2} \phi_{\lambda}^{L}\left(x_{2}\right)\right]\right|  \tag{191}\\
& +\left|g_{\lambda}\right|\left|\mathrm{d}_{\lambda}\left[w_{\lambda}^{1 / 2} \phi_{\lambda}^{L}\left(x_{1}\right)\right]\right|\left|\left[w_{\lambda}^{1 / 2} \phi_{\lambda}^{L}\left(x_{2}\right)\right]\right|  \tag{192}\\
& +\left|g_{\lambda}\right|\left|\left[w_{\lambda}^{1 / 2} \phi_{\lambda}^{L}\left(x_{1}\right)\right]\right|\left|\mathrm{d}_{\lambda}\left[w_{\lambda}^{1 / 2} \phi_{\lambda}^{L}\left(x_{2}\right)\right]\right|  \tag{193}\\
\left|\mathrm{d}_{\lambda} F\left(\lambda, x_{1}, x_{2}\right)\right|=\mid & \mathrm{d}_{\lambda}\left\{g_{\lambda}\left[w_{\lambda}^{1 / 2} \phi_{\lambda}^{L}\left(x_{1}\right)\right]\left[w_{\lambda}^{1 / 2} \phi_{\lambda}^{L}\left(x_{2}\right)\right]\right\} \mid  \tag{194}\\
<\mid & \left|\mathrm{d}_{\lambda} g_{\lambda}\right|(3)(1+4 q)\left(x_{1}+\kappa\right)^{2} \exp \left(-16^{-1} \epsilon_{1} \lambda\right)  \tag{195}\\
& \quad \times(3)(1+4 q)\left(x_{2}+\kappa\right)^{2} \exp \left(-16^{-1} \epsilon_{1} \lambda\right)  \tag{196}\\
& +\widehat{g}_{0}(q)(6)(1+4 q)\left(x_{1}+\kappa\right)^{2} \exp \left(-16^{-1} \epsilon_{1} \lambda\right)  \tag{197}\\
& \times(3)(1+4 q)\left(x_{2}+\kappa\right)^{2} \exp \left(-16^{-1} \epsilon_{2} \lambda\right)  \tag{198}\\
& +\widehat{g}_{0}(q)(3)(1+4 q)\left(x_{1}+\kappa\right)^{2} \exp \left(-16^{-1} \epsilon_{1} \lambda\right)  \tag{199}\\
& \times(6)(1+4 q)\left(x_{2}+\kappa\right)^{2} \exp \left(-16^{-1} \epsilon_{2} \lambda\right)  \tag{200}\\
\leq & c_{1}(q, \delta)\left(x_{1}+\kappa\right)^{2}\left(x_{2}+\kappa\right)^{2} \exp \left[-16^{-1}\left(\epsilon_{1}+\epsilon_{2}\right) \lambda\right]\left|\mathrm{d}_{\lambda} g_{\lambda}\right|  \tag{201}\\
= & \mathcal{O}\left(\lambda^{-1} \log ^{-1} \lambda\right) \tag{202}
\end{align*}
$$

in the norm of $\mathscr{B}$, where $c_{1}(q, \delta)$ is a constant.
Property (5): For $\lambda \searrow 0$ one has

$$
\begin{align*}
\left|\mathrm{d}_{\lambda}^{2} F\left(\lambda, x_{1}, x_{2}\right)\right| & =\left|\mathrm{d}_{\lambda}^{2}\left\{g_{\lambda}\left[w_{\lambda}^{1 / 2} \phi_{\lambda}^{L}\left(x_{1}\right)\right]\left[w_{\lambda}^{1 / 2} \phi_{\lambda}^{L}\left(x_{2}\right)\right]\right\}\right|  \tag{203}\\
& \leq c_{2}(q, \delta)\left(x_{1}+\kappa\right)^{3}\left(x_{2}+\kappa\right)^{3} \exp \left[-16^{-1}\left(\epsilon_{1}+\epsilon_{2}\right) \lambda\right]\left|\mathrm{d}_{\lambda}^{2} g_{\lambda}\right|  \tag{204}\\
& =\mathcal{O}\left(\lambda^{-2} \log ^{-2} \lambda\right) \tag{205}
\end{align*}
$$

in the norm of $\mathscr{B}$, where $c_{2}(q, \delta)$ is a constant.

## CHAPTER 2

# The nonlinear stationary problem 

## 1. Results

Consider the discrete NLS

$$
\begin{equation*}
i \partial_{t} w=L_{0} w-|w|^{2 \sigma} w, \quad 1 \leq \sigma \in \mathbb{Z} \tag{206}
\end{equation*}
$$

where $w: \mathbb{R}_{t} \times \mathbb{Z}_{+} \rightarrow \mathbb{C}$. The existence of a $u: \mathbb{Z}_{+} \rightarrow \mathbb{C}$ which satisfies the nonlinear finite difference equation

$$
\begin{equation*}
L_{0} u=\zeta u+|u|^{2 \sigma} u \tag{207}
\end{equation*}
$$

furnishes a stationary state of the discrete NLS of the form $w(t)=e^{-i \zeta t} u$. One expects that, due to the attractive nature of the nonlinearity, a negative "nonlinear eigenvalue", $\zeta=-a<0$, will allow the existence of a sharply peaked, monotonically decaying "ground state soliton". We will therefore exclusively look for solutions to

$$
\begin{equation*}
L_{0} u=-a u+u^{2 \sigma+1} \tag{208}
\end{equation*}
$$

where $u: \mathbb{Z}_{+} \rightarrow \mathbb{R}_{+}$and $a>0$. Solutions with these characteristics are self-focusing and tend to be sharply localized. They are therefore termed solitary waves or solitons generally.

Theorem 4. There exists a $\mu_{*}>0$ such that for each $\mu>\mu_{*}$ there exists a solution to Equation (208) with $\zeta=-\mu<0$ and $u=\alpha_{\mu}$, which is:
(1) positive: $\alpha_{\mu}(x)>0$ for all $x \in \mathbb{Z}_{+}$
(2) monotonically decaying: $\alpha_{\mu}(x+1)-\alpha_{\mu}(x)<0$ for all $x \in \mathbb{Z}_{+}$
(3) absolutely integrable: $\alpha_{\mu} \in \ell^{1}$

Definition 1.1. We define $\Pi:=I-\Pi_{0}$ and write $\widehat{v} \equiv \Pi v$ for all $v \in \mathscr{H}$ and write $\widehat{A} \equiv \Pi A$ for all $A \in \mathcal{B}(\mathscr{H})$.

The proof of Theorem (4) will proceed as follows:
(1) Consider Equation (208) with $\zeta=-a<0$, where $a$ is a free parameter, $u: \mathbb{Z}_{+} \rightarrow \mathbb{R}_{+}$, and $p:=2 \sigma+1: L_{0} u=-a u+u^{p}$. Split this equation into a boundary piece and a tail piece by applying $\Pi_{0}$ and $\Pi$ respectively.
(2) We take $b:=u(0)$ to be a fixed constant and iterate the tail piece of Equation (208) via

$$
\begin{equation*}
u_{n+1}(a, b)=\widehat{\psi}_{-a} b^{p}+\widehat{R}_{-a}^{L_{0}} u_{n}^{p}(a, b) \tag{209}
\end{equation*}
$$

$u_{n}(a, b) \in \mathscr{H}$ for all $n$. We show that for large enough $a$ this iteration converges pointwise monotonically and that $\left\|u_{n}(a, b)\right\|_{1} \leq s_{n}(a)$ is bounded as $n \nearrow \infty$, for a sequence of constants $\left\{s_{n}(a)\right\}_{n=0}^{\infty}$. We define the limit of this iteration to be $u_{*}(a, b) \equiv \lim _{n} \nearrow_{\infty} u_{n}(a, b)$.
(3) The construction of $u_{*}(a, b)$ sets $u(1)=u_{*}(a, b ; 1) \equiv q(a, b)$. We substitute this into the boundary piece of Equation (208) which then takes the form

$$
\begin{equation*}
0=b^{p}-(a+1) b+q(a, b) . \tag{210}
\end{equation*}
$$

We will now take $b$ to be a variable. If the solution, $u$, is positive and monotonically decaying then one must have

$$
\begin{equation*}
0<u(1)<u(0) \quad \Rightarrow \quad a^{(p-1)^{-1}}<b<(a+1)^{(p-1)^{-1}} \tag{211}
\end{equation*}
$$

We show that for all $a$ sufficiently large there is a unique $b=b_{*}(a) \in$ $\left(a^{(p-1)^{-1}},(a+1)^{(p-1)^{-1}}\right)$ which solves the boundary equation.
(4) We define the solution we desire, $\alpha_{\mu}$, by

$$
\alpha_{\mu}(x):=\left\{\begin{array}{cl}
b_{*}(\mu) & , x=0  \tag{212}\\
u_{*}\left(\mu, b_{*}(\mu) ; x\right) & , x>0
\end{array} .\right.
$$

(5) The three properties (i.e. positivity, monotonicity, $\ell^{1}$ ) of the solution will then be verified in turn.

Typically one can arrive at the existence of a soliton with such properties via variational or rearrangement arguments. We will use much more elementary techniques which yield yet other properties due to the dependence on explicit constructions. One such result which will be of use later on is

## Proposition 2.

$$
\begin{equation*}
\left\|\widehat{\alpha}_{\mu}\right\|_{1} \leq s(\mu)=\mu^{-1}(\mu+1)^{(p-1)^{-1}}+\mu s^{p}(\mu)=\mu^{-(p-1)^{-1}(p-2)}+\mathcal{O}\left(\mu^{-(p-1)^{-1}(2 p-3)}\right) . \tag{213}
\end{equation*}
$$

## 2. Existence of $\alpha_{\mu}$

2.1. Away from the boundary. Consider two forms of Equation (208):

$$
\begin{align*}
L_{0} u & =-a u+u^{p}  \tag{I}\\
u & =R_{-a}^{L_{0}} u^{p} . \tag{II}
\end{align*}
$$

One may project the equation of form (II) to a "tail" piece by applying $\Pi$ :

$$
\begin{equation*}
u=R_{-a}^{L_{0}}\left(\Pi_{0}+\Pi\right) u^{p} \quad \Rightarrow \quad \widehat{u}=\widehat{\psi}_{-a} u^{p}(0)+\widehat{R}_{-a}^{L_{0}} \widehat{u}^{p} \tag{214}
\end{equation*}
$$

We will fix $u(0) \equiv b$ and iterate by substituting the LHS into the RHS. We will show the conditions under which this converges and

Definition 2.1. Let $u(0) \equiv b$ be a fixed parameter which satisfies $a^{(p-1)^{-1}}<b<$ $(a+1)^{(p-1)^{-1}}$. Let $\left\{u_{n}(a, b)\right\}_{n=0}^{\infty}$ be a sequence of vectors in $\mathscr{H}$ defined by a fixed $u_{0}(a, b)$ and inductively by

$$
\begin{equation*}
u_{n+1}(a, b)=\widehat{\psi}_{-a} b^{p}+\widehat{R}_{-a}^{L_{0}} u_{n}^{p}(a, b) \tag{215}
\end{equation*}
$$

such that $u_{n}(a, b)=\widehat{u}_{n}(a, b)$.

The requirement that $a^{(p-1)^{-1}}<b<(a+1)^{(p-1)^{-1}}$ follows from $0<u(1)<u(0)$ such that $u(0)=b$ and $u(x)$ in monotonically decreasing for increasing $x$.

Proof. $\psi_{-a}(x)>0$ for all $a \in \mathbb{R}_{+}$and $x \in \mathbb{Z}_{+}$. Therefore

$$
\begin{align*}
\left\|\psi_{-a}\right\|_{1} & =\sum_{x=0}^{\infty} \psi_{-a}(x)  \tag{216}\\
& =\sum_{x=0}^{\infty} e^{a} \sum_{k=0}^{x}(-1)^{k}\binom{x}{k} E_{k+1}(a)=\sum_{x=0}^{\infty} e^{a} \int_{1}^{\infty} \mathrm{d} t e^{-a t} t^{-1}\left(1-t^{-1}\right)^{x}  \tag{217}\\
& =e^{a} \int_{1}^{\infty} \mathrm{d} t e^{-a t} t^{-1} \sum_{x=0}^{\infty}\left(1-t^{-1}\right)^{x}=\int_{0}^{\infty} \mathrm{d} t e^{-a t}=a^{-1} \tag{218}
\end{align*}
$$

Lemma 2.2. Let $\left\{s_{n}(a)\right\}_{n=0}^{\infty}$ be a sequence of nonnegative real numbers defined by a fixed $s_{0}(a)$ and inductively by

$$
\begin{equation*}
s_{n+1}(a)=a^{-1} r(a)+a^{-1} s_{n}^{p}(a), \tag{219}
\end{equation*}
$$

where

$$
\begin{equation*}
r(a):=(a+1)^{(p-1)^{-1}} . \tag{220}
\end{equation*}
$$

If $\left\|u_{j}(a, b)\right\|_{1} \leq s_{j}(a)$ for some $j$ then $\left\|u_{k}(a, b)\right\|_{1}<s_{k}(a)$ for all $j<k$.

Proof. Consider the known bound [1]

$$
\begin{equation*}
(z+n)^{-1}<e^{z} E_{n}(z) \leq(z+n-1)^{-1}, \quad 0<z \in \mathbb{R} \tag{221}
\end{equation*}
$$

Since $\psi_{-a}(0)=e^{a} E_{1}(a)$, one then has that

$$
\begin{equation*}
\left\|\widehat{\psi}_{-a}\right\|_{1}=\left\|\psi_{-a}\right\|_{1}-\psi_{-a}(0)<a^{-1}(a+1)^{-1} \tag{222}
\end{equation*}
$$

Then, since $b<(a+1)^{(p-1)^{-1}}$ one has

$$
\begin{equation*}
\left\|\widehat{\psi}_{-a} b^{p}\right\|_{1}<a^{-1} r(a) \tag{223}
\end{equation*}
$$

One may then observe that

$$
\begin{equation*}
\left\|u_{j+1}\right\|_{1}<a^{-1} r(a)+a^{-1}\left\|u_{j}^{p}\right\|_{1} \leq a^{-1} r(a)+a^{-1} s_{j}^{p}(a)=s_{j+1}(a) . \tag{224}
\end{equation*}
$$

Lemma 2.3. Let $g: \mathbb{R}_{+}^{2} \rightarrow \mathbb{R}$ be specified by

$$
\begin{equation*}
g(a, s):=a^{-1} r(a)+a^{-1} s^{p}-s . \tag{225}
\end{equation*}
$$

and $s_{\min }(a):=(a / p)^{(p-1)^{-1}}$. For sufficiently large $a>0$ it is the case that $g(a, s)$ has exactly two roots in $s: s_{-}(a)$ which satisfies $0<s_{-}(a)<s_{\min }(a)$ and $s_{+}(a)$ which satisfies $s_{\min }(a)<s_{+}(a)$.

Proof. One may observe that $g(a, s)$ has a global minimum at $s=s_{\text {min }}(a)$. It is the case that

$$
\begin{align*}
& g\left(a, s_{\min }(a)\right)=a^{-1} r(a)+a^{-1}(a / p)^{p(p-1)^{-1}}-a^{(p-1)^{-1}} p^{-(p-1)^{-1}}  \tag{226}\\
& =a^{-1}(a+1)^{(p-1)^{-1}}-a^{(p-1)^{-1}} p^{-(p-1)^{-1}}\left(1-p^{-1}\right)  \tag{227}\\
& \mathrm{d}_{a} g\left(a, s_{\min }(a)\right)=-a^{-1}(a+1)^{(p-1)^{-1}}\left[a^{-1}-(a+1)^{-1}(p-1)^{-1}\right]  \tag{228}\\
& \\
& \quad \quad-a^{-(p-2)(p-1)^{-1}} p^{-(p-1)^{-1}}(p-1)^{-1}\left(1-p^{-1}\right)<0, \quad \forall a>0
\end{align*}
$$

since

$$
\begin{equation*}
a^{-1}-(a+1)^{-1}(p-1)^{-1}>0, \quad \forall a>0 . \tag{230}
\end{equation*}
$$

One has that

$$
\begin{equation*}
\lim _{a \nearrow \infty} g\left(a, s_{\min }(a)\right)=-\infty \tag{231}
\end{equation*}
$$

and

$$
\begin{align*}
g(a, 0) & =a^{-1} r(a)>0, \quad \forall a>0  \tag{232}\\
\lim _{a \nearrow \infty} g(a, 0) & =\infty  \tag{233}\\
\lim _{s \nearrow \infty} g(a, s) & =\infty, \quad \forall a>0 \tag{234}
\end{align*}
$$

By intermediate value theorem there must be at least one root. By Descartes rule of signs, $g(a, s)$ has either 0 or 2 positive roots in $s$. Therefore $g(a, s)$ has exactly 2 roots for all sufficiently large $a>0$.

Definition 2.2. Let $a_{0}>0$ be the unique value of a such that $g(a, s)$ has two distinct roots in $s$ for all $a>a_{0}$.
$a_{0}>0$ is the unique value of $a$ for which that $g(a, s)$ has one root in $s$ of multiplicity 2.

Lemma 2.4. Let $h: \mathbb{R}_{+}^{2} \rightarrow \mathbb{R}_{+}$be defined by

$$
\begin{equation*}
h(a, s):=a^{-1} r(a)+a^{-1} s^{p} . \tag{235}
\end{equation*}
$$

The map $h(a, \cdot): s \mapsto h(a, s)$ is contractive on the domain $0 \leq s<s_{\min }(a)$ for all $a>a_{0}$.

## Proof.

$$
\begin{align*}
\left.\partial_{s} g(a, s)\right|_{s=0} & =-1  \tag{236}\\
\left.\partial_{s} g(a, s)\right|_{s=s_{\min }(a)} & =0  \tag{237}\\
\partial_{s}^{2} g(a, s) & >0, \quad \forall s>0  \tag{238}\\
\Rightarrow \quad 0 \leq \partial_{s} h(a, s) & <1, \quad \forall s: 0 \leq s<s_{\min }(a) \tag{239}
\end{align*}
$$

By mean value theorem, for any $s, s_{1} \in\left[0, s_{\min }\right)$ there exists an $s_{2}<\left|s-s_{1}\right|$ such that

$$
\begin{equation*}
\frac{h(a, s)-h\left(a, s_{1}\right)}{s-s_{1}}=\left.\partial_{s} h(a, s)\right|_{s=s_{2}}<1 \tag{240}
\end{equation*}
$$

which completes the proof.

Lemma 2.5. For all $a>a_{0}$ and for sufficiently small $\left\|u_{0}(a, b)\right\|_{1} \geq 0$ one has that $\lim _{n \neq \infty}\left\|u_{n}(a, b)\right\|_{1}<\infty$.

Proof. Let $a>a_{0}$. Given the iteration $s_{n+1}(a)=h\left(a, s_{n}(a)\right)$, the choice of any $s_{0}(a)$ which satisfies $0 \leq s_{0}(a) \leq s_{-}(a)$ gives a sequence $\left\{s_{n}(a)\right\}_{n=0}^{\infty}$ which converges to $\lim _{n \nearrow \infty} s_{n}(a)=s_{-}(a)<\infty$ monotonically from below as $n \nearrow \infty$ for all $a>a_{0}$. For $0 \leq\left\|u_{0}(a, b)\right\|_{1} \leq s_{0}(a)$ it is the case that $\lim _{n \nmid \infty}\left\|u_{n}(a, b)\right\|_{1} \leq \lim _{n \nmid \infty} s_{n}(a)=$ $s_{-}(a)<\infty$.

LEmma 2.6. Let $s_{*}(a):=\lim _{n \nmid \infty} s_{n}(a)=s_{-}(a)$ for all $a>a_{0}$. One has that $s_{*}(a) \searrow 0$ monotonically as a $\nearrow \infty$ for all $a>a_{0}$.

Proof. Let $a>a_{0}$. Consider the graph of the map $s \mapsto h(a, s)=a^{-1} r(a)+a^{-1} s^{p}$. It is the case that $\partial_{a} h(a, s)<0$ and $h(a, s) \geq s$ for $0 \leq s \leq s_{*}(a)$, with equality for $h\left(a, s_{*}(a)\right)=s_{*}(a)$, which completes the proof.

Lemma 2.7. We define

$$
\begin{align*}
\theta_{1}(a, b):= & \widehat{\psi}_{-a} b^{p}  \tag{241}\\
\theta_{k}(a, b):= & \widehat{R}_{-a}^{L_{0}} \sum_{m_{1}=1}^{p}\binom{p}{m_{1}} \theta_{1}^{p-m_{1}}(a, b) \cdots \sum_{m_{i}=1}^{m_{j-1}}\binom{m_{j-1}}{m_{j}} \theta_{j}^{m_{j}-m_{j-1}}(a, b) \cdots  \tag{242}\\
& \times \sum_{m_{k-2}=1}^{m_{k-3}}\binom{m_{k-3}}{m_{k-2}} \theta_{k-2}^{m_{k-3}-m_{k-2}}(a, b) \theta_{k-1}^{m_{k-2}}(a, b) . \tag{243}
\end{align*}
$$

If $u_{0}(a, b)=0$ then one has

$$
\begin{equation*}
u_{n}(a, b)=\sum_{k=1}^{n} \theta_{k}(a, b) \tag{244}
\end{equation*}
$$

Proof. One may verify by induction:

$$
\begin{equation*}
u_{1}(a, b)=\theta_{1}(a, b) \tag{245}
\end{equation*}
$$

$u_{n+1}(a, b)=\theta_{1}(a, b)+\widehat{R}_{-a}^{L_{0}} u_{n}^{p}(a, b)$

$$
\begin{equation*}
=\theta_{1}(a, b)+\widehat{R}_{-a}^{L_{0}}\left[\sum_{k=1}^{n} \theta_{k}(a, b)\right]^{p} \tag{247}
\end{equation*}
$$

$$
=\theta_{1}(a, b)+\widehat{R}_{-a}^{L_{0}}\left\{\theta_{1}^{p}(a, b)+\sum_{m_{1}=1}^{p}\binom{p}{m_{1}} \theta_{1}^{p-m_{1}}(a, b)\left[\sum_{k=2}^{n} \theta_{k}(a, b)\right]^{m_{1}}\right\}
$$

$$
=\theta_{1}(a, b)+\widehat{R}_{-a}^{L_{0}}\left[\theta_{1}^{p}(a, b)+\sum_{m_{1}=1}^{p}\binom{p}{m_{1}} \theta_{1}^{p-m_{1}}(a, b) \theta_{2}^{m_{1}}(a, b)+\cdots\right.
$$

$$
+\sum_{m_{1}=1}^{p}\binom{p}{m_{1}} \theta_{1}^{p-m_{1}}(a, b) \cdots \sum_{m_{j}=1}^{m_{j-1}}\binom{m_{j-1}}{m_{j}} \theta_{j}^{m_{j}-m_{j-1}}(a, b) \cdots
$$

$$
\times \sum_{m_{k-2}=1}^{m_{k-3}}\binom{m_{k-3}}{m_{k-2}} \theta_{k-2}^{m_{k-3}-m_{k-2}}(a, b) \theta_{k-1}^{m_{k-2}}(a, b)
$$

$$
+\sum_{m_{1}=1}^{p}\binom{p}{m_{1}} \theta_{1}^{p-m_{1}}(a, b) \cdots \sum_{m_{j}=1}^{m_{j-1}}\binom{m_{j-1}}{m_{j}} \theta_{j}^{m_{i}-m_{j-1}}(a, b) \cdots
$$

$$
\left.\times \sum_{m_{n-1}=1}^{m_{n-2}}\binom{m_{n-2}}{m_{n-1}} \theta_{n-1}^{m_{n-2}-m_{n-1}}(a, b) \theta_{n}^{m_{n-1}}(a, b)\right]
$$

$$
\begin{equation*}
=\theta_{1}(a, b)+\theta_{2}(a, b)+\cdots+\theta_{j}(a, b)+\cdots+\theta_{n+1}(a, b) \tag{254}
\end{equation*}
$$

$$
\begin{equation*}
=\sum_{k=1}^{n+1} \theta_{k}(a, b) \tag{255}
\end{equation*}
$$

Lemma 2.8. For $u_{0}(a, b)=0$ one has that $u_{n+1}(a, b ; x)>u_{n}(a, b ; x)$ for all $n \geq$ $0, x>0$ (strict pointwise increase in $n$ for all $x>0$ ) for all $a>a_{0}$.

Proof. Since $\theta_{n}(a, b ; x)>0, \forall x>0$, and $u_{n+1}(a, b)-u_{n}(a, b)=\theta_{n+1}(a, b)$ it must be the case that the sequence $\left\{u_{n}(a, b ; x)\right\}_{n=0}^{\infty}$ is strictly increasing in $n$ for all $x>0$ and for all $a>a_{0}$.

Lemma 2.9. For $u_{0}(a, b)=0$ and all $a>a_{0}$ one has that $\lim _{n \neq \infty} u_{n}(a, b)=$ $\sum_{n=1}^{\infty} \theta_{n}(a, b) \in \ell^{1}$.

Proof. Let $a>a_{0}$. One may observe that that $\lim _{n \nmid \infty}\left\|u_{n}(a, b)\right\|_{1} \leq s_{*}(a)<$ $\infty \Rightarrow \lim _{n \nearrow \infty} u_{n}(a, b ; x)<s_{*}(a)<\infty$ for all $x \in \mathbb{Z}_{+}$. Then by monotonic increase of $u_{n}(a, b ; x)$ in $n$ for all $x>0$, it follows that $u_{n}(a, b ; x)=\sum_{k=1}^{n} \theta_{k}(a, b ; x)$ exists for each $n, x$ and uniquely determines $\lim _{n \nearrow_{\infty}} u_{n}(a, b)$.

Lemma 2.10. Let $u_{*}(a, b):=\sum_{n=1}^{\infty} \theta_{n}(a, b)=\lim _{n} \nmid \infty u_{n}(a, b)$ for all $a>a_{0}$. One has that $u_{*}(a, b ; x)$ is a monotonically increasing function in $b \geq 0$ for all $x>0$ and for all $a>a_{0}$.

Proof. $u_{*}(a, b ; x):=\sum_{n=1}^{\infty} \theta_{n}(a, b ; x)$ can be represented as a power series in $b$ for all $x>0$ with only positive powers and positive coefficients.
2.2. At the boundary. We now consider the equation of form (I). We can project it onto a "boundary" piece by applying $\Pi_{0}$ :

$$
\begin{equation*}
0=\left(\Pi_{0}+\Pi\right)\left(-L_{0} u-a u+u^{p}\right) \quad \Rightarrow \quad 0=u^{p}(0)-(a+1) u(0)+u(1) \tag{256}
\end{equation*}
$$

Given $u_{*}(a, b)=\sum_{k=1}^{\infty} \theta_{k}(a, b)$ we will substitute $b=u(0)$ and $q(a, b) \equiv u_{*}(a, b ; 1)=$ $u(1)$ in the boundary equation and thereby consider

$$
\begin{equation*}
0=b^{p}-(a+1) b+q(a, b) \tag{257}
\end{equation*}
$$

on the interval $a^{(p-1)^{-1}} \leq b \leq(a+1)^{(p-1)^{-1}}$.


Figure 1. The dashed line is the graph of $f_{+}(a, b)$, the solid is of $f(a, b)$, and the dotted is of $f_{-}(a, b)$.

## Definition 2.3.

$$
\begin{align*}
& b_{-}(a):=a^{(p-1)^{-1}}, \quad b_{+}(a):=(a+1)^{(p-1)^{-1}}  \tag{258}\\
& \Sigma_{a}:=\left\{b \in \mathbb{R}: b_{-}(a) \leq b \leq b_{+}(a)\right\}  \tag{259}\\
& q_{-}(a, b):=0, \quad q(a, b):=u_{*}(a, b ; 1), \quad q_{+}(a, b):=b  \tag{260}\\
& f(a, b, q):=b^{p}-(a+1) b+q  \tag{261}\\
& f_{-}(a, b):=f\left(a, b, q_{-}(a, b)\right), \quad f_{*}(a, b):=f(a, b, q(a, b)),  \tag{262}\\
& \quad f_{+}(a, b):=f\left(a, b, q_{+}(a, b)\right) \tag{263}
\end{align*}
$$

Lemma 2.11. $q(a, b)>q_{-}(a, b)$ for all $b \in \Sigma_{a}$ and for all $a>a_{0}$.

Proof. Let $a>a_{0}$.

$$
\begin{align*}
q(a, b) & =\sum_{n=1}^{\infty} \theta_{n}(a, b ; 1)>\theta_{1}(a, b ; 1)=\widehat{\psi}_{-a}(1) b^{p}  \tag{264}\\
& =e^{a}\left[E_{1}(a)-E_{2}(a)\right] b^{p}>0=q_{-}(a, b) \tag{265}
\end{align*}
$$

since $E_{2}(a)<E_{1}(a), \forall a>0[1]$.

Lemma 2.12. There exists an $a_{1} \geq a_{0}$ such that $q(a, b)<q_{+}(a, b)$ for all $b \in \Sigma_{a}$ and for all $a>a_{1}$.

Proof. Let $a>a_{0}$. It is the case that $\lim _{a \not ~_{\infty}} s_{*}(a)=0$ monotonically and $\lim _{a \rightarrow \infty} b_{-}(a)=\infty$ monotonically so there exists an $a_{1} \geq a_{0}$ such that $s_{*}(a)<$ $b_{-}(0), \forall a>a_{1}$. Clearly $q(a, b)<s_{*}(a)$. Then $q(a, b)<s_{*}(a)<b_{-}(a) \leq b=$ $q_{+}(a, b), \forall a>a_{1}$ and $\forall b \in \Sigma_{a}$.

LEmma 2.13. Let $a_{1}>0$ be the smallest value for which $a_{1} \geq a_{0}$ and for which $q(a, b)<q_{+}(a, b)$ for all $b \in \Sigma_{a}$ and for all $a>a_{1}$. Define $a_{2}:=(p-1)^{-1}, a_{3}:=$ $\max \left\{a_{1}, a_{2}\right\}$. One has that:

- $f_{-}(a, b)$ is negative and monotonically increasing for all $b \in \Sigma_{a} \backslash b_{+}(a)$ and for all $a>a_{3}$.
- $f_{+}(a, b)$ is positive and monotonically increasing in $b$ for all $b \in \Sigma_{a} \backslash b_{-}(a)$ and for all $a>a_{3}$.

Proof. One may observe that $f_{-}(a, b)<f_{*}(a, b)<f_{+}(a, b)$ for all $a>a_{1}$ and all $b \in \Sigma_{a} . f_{-}\left(a, b_{+}(a)\right)=0$ and $f_{+}\left(a, b_{-}(a)\right)=0$ for all $a \geq 0$.

For $f_{-}(a, b)$ : Let $a>a_{3}$. It follows that

$$
\begin{equation*}
\partial_{b} f_{-}(a, b)=p b^{p-1}-(a+1)>p b_{-}^{p-1}(a)-(a+1)=a(p-1)-1 . \tag{266}
\end{equation*}
$$

Then $\partial_{b} f_{-}(a, b)>0, \forall b \in \Sigma_{a}$ and for all $a>a_{3}$. Since $f_{-}\left(a, b_{+}(a)\right)=0$ one has that $f_{-}(a, b)<0$ for all $b \in \Sigma_{a} \backslash b_{+}(a)$.

For $f_{+}(a, b)$ : Let $a>a_{3}$.

$$
\begin{equation*}
\partial_{b} f_{+}(a, b)=p b^{p-1}-a \geq p b_{-}^{p-1}(a)-a=a(p-1)>0, \forall a>0 . \tag{267}
\end{equation*}
$$

Then since $f_{+}\left(a, b_{-}(a)\right)=0$ one finds that $f_{+}(a, b)>0 \forall b \in \Sigma_{a} \backslash b_{-}(a)$.

Lemma 2.14. $f_{*}(a, b)$ is monotonically increasing in $b$ and has exactly one root in $b$ for all $b \in \Sigma_{a}$ and for all $a>a_{3}$.

Proof. Let $a>a_{3} . f_{-}(a, b)$ is monotonically increasing in $b$ for all $b \in \Sigma_{a} . q(a, b)$ is monotonically increasing in $b$ for all $b>0 . f_{*}(a, b)=f_{-}(a, b)+q(a, b)$ so $f_{*}(a, b)$ must be monotonically increasing in $b$ for all $b \in \Sigma_{a}$.

Since $f_{*}\left(a, b_{-}(a)\right)<0=f_{+}\left(b_{-}(a)\right)=f_{-}\left(a, b_{+}(a)\right)<f_{*}\left(a, b_{+}(a)\right)$, by intermediate value theorem there must be at least one $b=b_{*}(a)$ for which $f_{*}(a, b)=0$ and since $f_{*}(a, b)$ is monotonically increasing on this interval, there must be exactly one such $b=b_{*}(a)$.

DEFINITION 2.4. Let $b_{*}(a)$ be the unique value of $b \in \Sigma_{a}$ for which $f_{*}(a, b)=0$ for all $a>a_{3}$ and define $\alpha_{a} \in \mathscr{H}$ via

$$
\alpha_{a}(x):=\left\{\begin{array}{cc}
b_{*}(a) & , x=0  \tag{268}\\
u_{*}\left(a, b_{*}(a)\right) & , x>0
\end{array},\right.
$$

for all $a>a_{3}$.

## 3. Properties of $\alpha_{\mu}$

### 3.1. Monotonicity. Consider

$$
\begin{equation*}
-L_{0} u=V(u) \tag{269}
\end{equation*}
$$

where $V(\cdot): \mathbb{R} \rightarrow \mathbb{R}$ satisfies $V(r)=0, V\left(r_{0}\right)=0$, and $V(r)>0$ for all $r \in\left(0, r_{0}\right)$ where $r_{0}>0$. It was shown in [10] that for any solution of this equation for which $\lim _{x \nearrow \infty} u(x)=0$, there exists an $x_{*}$ such that $u(x)$ is a monotonically decaying for increasing $x \geq x_{*}$. We will apply their argument, section 3 of [10], to show that for sufficiently large $a$ one must have that $\alpha_{a}(x)$ is monotonically decaying in $x$ for all $x \in \mathbb{Z}_{+}$. For our case we have that $V(r)=a r-r^{p}, r \in \mathbb{R}$, which satisfies the desired criteria.

Equation (269) may be summed to give an alternative finite difference equation. One has the two forms:

$$
\begin{align*}
(x+1)[u(x+1)-u(x)]-x[u(x)-u(x-1)] & =V(u(x)),  \tag{270}\\
u(x+1)-u(x) & =(x+1)^{-1} \sum_{y=0}^{x} V\left(u_{y}\right) . \tag{271}
\end{align*}
$$

The argument proceeds as follows.

Lemma 3.1. There exists an $a_{4} \geq a_{3}$ such that $s_{*}(a)<b_{-}(a)$ for all $a>a_{4}$.

Proof. Let $a>a_{3}$. It is the case that $s_{*}(a) \searrow 0$ monotonically as $a \nearrow \infty$. It is clear that $b_{-}(a) \nearrow \infty$ monotonically as $a \nearrow \infty$.

Lemma 3.2. Let $u \in \mathscr{H}$ solve Equation (269). Let $c_{-}$and $c_{+}$be constants which satisfy $0<c_{-}<u(0)$ and $0<\|\Pi u\|<c_{+}<\infty$. If additionally $c_{+}<c_{-}$then $u(x)$ is monotonically decreasing in $x$ for all $x \in \mathbb{Z}_{+}$.

Proof. Let $u$ solve Equation (269). Let $c_{-}$and $c_{+}$be constants which satisfy $0<c_{-}<u(0)$ and $0<\|\Pi u\|<c_{+}<\infty$. Furthermore let $c_{+}<c_{-}$. One has that
$u(x)<u(0)$ for all $x>0$ and in particular $u(1)<u(0)$. Assume that there exists an $0<x_{0} \in \mathbb{Z}_{+}$such that $u\left(x_{0}+1\right)-u\left(x_{0}\right)=0$ for $u\left(x_{0}\right), u\left(x_{0}+1\right) \in\left(0, c_{-}\right)$. One then has $u\left(x_{0}+2\right)-u\left(x_{0}+1\right)>0$ since $\left(x_{0}+2\right)\left[u\left(x_{0}+2\right)-u\left(x_{0}+1\right)\right]=V\left(u\left(x_{0}+1\right)\right)>0$. Then assume generally that there exists an $0<x_{1} \in \mathbb{Z}_{+}$such that $u\left(x_{1}+1\right)-u\left(x_{1}\right)>0$. This gives that $\sum_{y=0}^{x_{1}} V(u(y))>0$ since $u\left(x_{1}+1\right)-u\left(x_{1}\right)=\left(x_{1}+1\right)^{-1} \sum_{y=0}^{x_{1}} V(u(y))>0$. One then has that

$$
\begin{align*}
& u\left(x_{1}+2\right)-u\left(x_{1}+1\right)=\left(x_{1}+2\right)^{-1} \sum_{y=0}^{x_{1}+1} V(u(y))  \tag{272}\\
&=\left(x_{1}+2\right)^{-1} V\left(u\left(x_{1}+1\right)\right)+\left(x_{1}+2\right)^{-1} \sum_{y=0}^{x_{1}} V(u(y))>0 \tag{273}
\end{align*}
$$

since $V\left(u\left(x_{1}+1\right)\right)$ and $\sum_{y=0}^{x_{1}} V(u(y))$ are both positive. If $u\left(x_{1}+2\right) \geq c_{-}$then one has a contradiction. On the other hand if $u\left(x_{1}+2\right)<c_{-}$then one may repeat the above process with $x_{1}$ replaced with $x_{1}+1$. Therefore if there exists an $x_{1}>0$ such that $u\left(x_{1}+1\right)-u\left(x_{1}\right)>0$, the $u(x)$ for the subsequent $x>x_{1}+1$ will continue to rise at least until $u\left(x_{3}\right) \geq c_{-}$, which is the point greater than which $V(x)$ remains negative, for some $x_{3}>x_{2}+1$. One therefore has a contradiction if $u(x)$ fails to be monotonically decreasing as $x \nearrow \infty$ for all $x \in \mathbb{Z}_{+}$.

LEmma 3.3. Let $\mu_{*}>0$ be the smallest value such that $\mu_{*} \geq a_{3}$ and such that $s_{*}(\mu)<b_{-}(\mu)$ for all $\mu>\mu_{*}$. One has that $\alpha_{\mu}(x)$ decreases monotonically as $x \nearrow \infty$ for all $x \in \mathbb{Z}_{+}$and for all $\mu>\mu_{*}$

Proof. Consider Lemma 4.1 and let $u=\alpha_{\mu}, c_{-}=b_{-}(\mu)$, and $c_{+}=s_{*}(\mu)$ for all $\mu>\mu_{*}$.
3.2. Asymptotic behavior. In [5] the asymptotic behavior of solutions of the finite difference nonsingular Sturm-Liouville problem $\Pi\left(L_{0}-\lambda\right) \Pi u=0$ were studied with various boundary conditions. One solution, $\phi_{\lambda}$, has the well-known asymptotic
behavior

$$
\begin{equation*}
\phi_{\lambda}(x) \sim e^{\lambda / 2} J_{0}(2 \sqrt{\lambda x}) \quad \text { as } \quad x \nearrow \infty, \tag{274}
\end{equation*}
$$

where $J_{0}(z)$ is the Bessel function of the first kind of degree 0 . They studied a particular solution, which they call $\Psi_{\lambda}$, which is a linear combination of $\phi_{\lambda}$ and $\psi_{\lambda}$. It satisfies the boundary conditions $\Psi_{\lambda}(0)=0$ and $\Psi_{\lambda}(1)=1 . \Psi_{\lambda}$ was shown to have the asymptotic behavior

$$
\begin{equation*}
\Psi_{\lambda}(x) \sim \pi e^{-\lambda / 2} Y_{0}(2 \sqrt{\lambda x})+e^{\lambda / 2} \mathcal{P} E_{1}(-\lambda) J_{0}(2 \sqrt{\lambda x}) \quad \text { as } \quad x \nearrow \infty \tag{275}
\end{equation*}
$$

where $Y_{0}(z)$ is the Bessel function of the second kind of degree 0 . One can determine the asymptotic behavior of $\psi_{-a}(x)$ as $x \nearrow \infty$ by finding the appropriate linear combination of Bessel functions such that upon analytic continuation, the real part is of the appropriate linear combination of the asymptotic forms of $\phi_{\lambda}(x)$ and $\Psi_{\lambda}(x)$ is monotonically decreasing as $x \nearrow \infty$. This combination must be monotonically decaying for $\lambda=-a<0$ one can straightforwardly determine that the asymptotic behavior must be of the form

$$
\begin{equation*}
\psi_{-a}(x) \sim 2 e^{a / 2} K_{0}(2 \sqrt{a x}) \sim e^{a / 2} \pi^{1 / 2}(a x)^{-1 / 4} e^{-2 \sqrt{a x}} \quad \text { as } \quad x \nearrow \infty \tag{276}
\end{equation*}
$$

where $K_{0}(z)$ is the modified Bessel function of the second kind of degree 0 .

LEmmA 3.4. $\phi_{-a}(x), \psi_{-a}(x)>0$ for all $x \in \mathbb{Z}_{+}$and $a>0 . \phi_{-a}(x)$ is monotonically increasing and $\psi_{-a}(x)$ is monotonically decreasing in increasing $x$ for all $a>0$.

Proof. One can see that $\phi_{-a}(x)$ is positive and monotonically increasing by inspection of $\phi_{-a}(x)=\sum_{k=0}^{x} \frac{a^{k}}{k!}\binom{x}{k}$. One can observe that $\psi_{-a}(x)$ is positive for all $x, a$

$$
\begin{equation*}
\psi_{-a}(x)=e^{a} \sum_{k=0}^{x}(-1)^{k}\binom{x}{k} \mathrm{E}_{k+1}(a)=e^{a} \int_{1}^{\infty} \mathrm{d} t e^{-a t} t^{-1}\left(1-t^{-1}\right)^{x}>0 \tag{277}
\end{equation*}
$$

as well as monotonically decreasing

$$
\begin{equation*}
\psi_{-a}(x+1)-\psi_{-a}(x)=-e^{a} \int_{1}^{\infty} \mathrm{d} t e^{-a t} t^{-2}\left(1-t^{-1}\right)^{x}<0 \tag{278}
\end{equation*}
$$

Lemma 3.5. The resolvent $R_{z}^{L_{0}}$ has the Sturm-Liouville (SL) representation

$$
R_{z}^{L_{0}}\left(x_{1}, x_{2}\right)= \begin{cases}\phi_{z}\left(x_{1}\right) \psi_{z}\left(x_{2}\right), & x_{1} \leq x_{2}  \tag{279}\\ \phi_{z}\left(x_{2}\right) \psi_{z}\left(x_{1}\right), & x_{1} \geq x_{2}\end{cases}
$$

Proof. The operator $L_{0}$ on $\mathscr{H}$ is a singular, second order, finite difference SturmLiouville operator. This is made manifest when put into SL form, $L_{0}=D_{+} M D_{-}$, where $D_{+}, D_{-}$are the respectively the usual forward and backward finite difference operators

$$
\begin{align*}
& D_{+} v(x)=v(x+1)-v(x)  \tag{280}\\
& D_{-} v(x)=\left\{\begin{array}{cl}
v(x)-v(x-1) & , x>0 \\
v(x) & , x=0
\end{array}\right. \tag{281}
\end{align*}
$$

and $M$ is lattice index multiplication operator $M v(x)=x v(x)$ for all $v \in \mathscr{T}$.
$L_{0}$ is singular at the boundary point $x=0$. When its domain and range are restricted to functions only on lattice points for $x>0$ it is the case that $L_{0}$ is a nonsingular operator. This restricted operator, $\Pi L_{0} \Pi$, is second order and nonsingular therefore $\Pi L_{0} \Pi u=z u, u \in \mathscr{T}$, admits two linearly independent solutions which satisfy linearly independent boundary conditions.

The finite difference Wronskian, also known as the Cassoratian, of two vectors $u, v \in \mathscr{T}$ is given by (see e.g. [17])

$$
\begin{equation*}
W[u, v](x)=u(x) v(x+1)-u(x+1) v(x) . \tag{282}
\end{equation*}
$$

A Jacobi operator $A \in \mathcal{L}(\mathscr{T})$ can be brought into the form

$$
\begin{equation*}
A v(x)=\eta(x) v(x+1)+\omega(x) v(x)+\eta(x-1) v(x-1), \quad \eta(x), \omega(x) \in \mathbb{R} \tag{283}
\end{equation*}
$$

As is specified by the finite difference Sturm-Liouville theory, if $A$ is a Sturm-Liouville operator, $u_{z}, v_{z}$ are linearly independent solutions of $A u=z u$, and $\eta(x) W\left[u_{z}, v_{z}\right](x)=$ $1 \forall x$ then

$$
R_{z}^{A}\left(x_{1}, x_{2}\right)= \begin{cases}u_{z}\left(x_{1}\right) v_{z}\left(x_{2}\right), & x_{1} \leq x_{2}  \tag{284}\\ u_{z}\left(x_{2}\right) v_{z}\left(x_{1}\right), & x_{1} \geq x_{2}\end{cases}
$$

This construction follows for $A=\Pi L_{0} \Pi, u_{z}=\Pi \phi_{z}, v_{z}=\Pi \psi_{z}$. Since $L_{0}$ is singular at $x=0$ one cannot adjust boundary conditions any further than fixing a scale factor for spectral solutions. One can simulate a second boundary condition with the introduction of either a linear perturbation or of an inhomogenous source supported at the singular point. This is to say that one can respectively consider the equations

$$
\begin{equation*}
\left(L_{0}-q \Pi_{0}\right) u=z u \quad \text { or } \quad L_{0} u=z u+q \chi_{0} \tag{285}
\end{equation*}
$$

where $q \in \mathbb{C}$ is a parameter, the tuning of which simulates the tuning of a second boundary condition. By taking the latter form with $q=1$ one may arrive at $\psi_{z}$ for the second solution. It therefore must be the case that $\Pi R_{z}^{L_{0}} \Pi=\Pi R_{z}^{\Pi L_{0} \Pi} \Pi$ with the above prescription and $R_{z}^{L_{0}}=R_{z}^{\Pi L_{0} \Pi}$ may be checked for $x=0$ directly at boundary values.

One may observe that one has $R_{-a}^{L_{0}}\left(x_{1}, x_{2}\right)>0$ for all $a>0$ and $0 \leq x_{1}, x_{2} \in \mathbb{Z}$.

Lemma 3.6. Consider the equation

$$
\begin{equation*}
\left(L_{0}+a\right) u=V(u) u \tag{286}
\end{equation*}
$$

where $V(\cdot): \mathbb{R} \rightarrow \mathbb{R}$ is continuous, locally bounded, and satisfies $\lim _{r \searrow_{0}} V(r) r=0$. If $u$ is a solution of this equation which is positive for all $x$ and for which $\lim _{x \gamma_{\infty}} u(x)=$

0 then $u(x) \sim c_{0} x^{-1 / 2} e^{-c_{1} \sqrt{x}}$ as $x \nearrow \infty$ for some $0<c_{0}, c_{1}<\infty$ and for each fixed $a>0$.

Proof. Consider $a>0$ a fixed constant. Let $\Pi_{\leq x_{*}}:=\sum_{x=0}^{x_{*}} \Pi_{x}$ and $\Pi_{>x_{*}}:=$ $I-\Pi_{\leq x_{*}}$ for some $x_{*} \in \mathbb{Z}_{+}$. Let $q:=\left\|\Pi_{>x_{*}} V(u)\right\|_{\text {op }}=\left\|\Pi_{>x_{*}} V(u)\right\|_{\infty}$. Furthermore let $x_{*}$ satisfy $0<q<a$ for all $a>0$. One may find

$$
\begin{align*}
\left(L_{0}+a\right) u & =V(u) u  \tag{287}\\
\Rightarrow \quad u & =\left[L_{0}+a-\Pi_{>x_{*}(a)} V(u)\right]^{-1} \Pi_{\leq x_{*}(a)} u  \tag{288}\\
& =R_{-a}^{L_{0}}\left[1-\Pi_{>x_{*}} V(u) R_{-a}^{L_{0}}\right]^{-1} \Pi_{\leq x_{*}} u  \tag{289}\\
& =R_{-a}^{L_{0}} \sum_{n=0}^{\infty}\left[\Pi_{>x_{*}} V(u) R_{-a}^{L_{0}}\right]^{n} \Pi_{\leq x_{*}} u \tag{290}
\end{align*}
$$

where the sum converges absolutely. One has from the Sturm-Liouville form of the resolvent

$$
\begin{equation*}
R_{-a}^{L_{0}} \chi_{x} \leq \phi_{-a}(x) R_{-a}^{L_{0}} \chi_{0}=\phi_{-a}(x) \psi_{-a} \tag{291}
\end{equation*}
$$

for all $a>0$ and all $x \in \mathbb{Z}_{+}$. One then has, for all $a>0$ and all $x \in \mathbb{Z}_{+}$

$$
\begin{align*}
u(x) & \leq R_{-a}^{L_{0}} \sum_{n=0}^{\infty}\left[q\left(u, x_{*}(a)\right) R_{-a}^{L_{0}}\right]^{n} \Pi_{\leq x_{*}(a)} u(x)  \tag{292}\\
& \leq\left[R_{-a+q}^{L_{0}} \sum_{y=0}^{x_{*}} u(y) \chi_{y}\right](x)  \tag{293}\\
& \leq \sum_{y=0}^{x_{*}} u(y) \phi_{-a+q}(y) \psi_{-a+q}(x) \tag{294}
\end{align*}
$$

One may therefore conclude

$$
\begin{equation*}
u(x) \leq c \psi_{-a+q}(x) \tag{295}
\end{equation*}
$$

where $c=\sum_{y=0}^{x_{*}} u(y) \phi_{-a+q}(y)<\infty$. One has that $\psi_{-a+q}(x) \sim c^{\prime} x^{-1 / 2} e^{-2 \sqrt{(a-q) x}}$ as $x \nearrow \infty$ with $a$ fixed, where $c^{\prime}$ is a constant that depends on $-a+q$ alone. By
inspection the appropriate constants $0<c_{0}, c_{1}<\infty$ may be found for each fixed $a>0$.

Lemma 3.7. $\alpha_{\mu}(x) \sim c_{0} x^{-1 / 2} e^{-c_{1} \sqrt{x}}$ as $x \nearrow \infty$ for all $x \in \mathbb{Z}_{+}$and for all $\mu>\mu_{*}$.

Proof. For $\mu>\mu_{*}$ it must be the case that $\alpha_{\mu}(x)$ exists as defined, be monotonically decreasing in $x$, and since it solves Equation (208), which is of the form given by Equation (269) it must have the desired asymptotic behavior.

Proof of Theorem (4). Existence and Property (1) are given by Definition (2.4) and arguments on which the definition depends. Property (2) is given by Lemma (3.3). Property (3) is given by Lemma (3.7).

Since we now have existence and the desired properties, we can estimate an upper bound on norms of the tail of $\alpha_{\mu}$. We may use the convergence for the iteration of the $s_{n}(\mu)$ functions for large values of $\mu$ to verify Proposition (2) directly.

## CHAPTER 3

## The linearized matrix problem

## 1. Results

We would like to study the evolution of solutions which, at least at an initial time $t=0$, are close to the stationary soliton solution $\widehat{u}(t)=e^{-i(-\mu t+\nu)} \alpha_{\mu}$, where $\nu \in \mathbb{R}$ is an arbitrary phase factor. We then consider the ansatz $u=e^{-i \theta}\left(\alpha_{\mu_{1}}+\beta\right)$ where

$$
\begin{gather*}
\theta(t):=-\int_{0}^{t} \mu_{1}(s) \mathrm{d} s+\nu_{1}(t)  \tag{296}\\
\mu_{1}(t)=\mu+\widehat{\mu}(t), \quad \mu_{1}(0)=\mu, \quad \nu_{1}(t)=\nu+\widehat{\nu}(t), \quad \nu_{1}(0)=\nu, \tag{297}
\end{gather*}
$$

and $\beta: \mathbb{R}_{t} \times \mathbb{Z}_{+} \rightarrow \mathbb{C}$ has the property that it and $\mathrm{d}_{t} \beta$ are small in norm at $t=0$. If one finds that $u(t) \rightarrow e^{-i\left(-\mu_{\infty} t+\nu_{\infty}\right)} \alpha_{\mu_{*}}$ in norm as $t \nearrow \infty$ for some $\mu_{\infty}, \nu_{\infty}$, for all $\mu, \nu$ and all sufficiently small $\beta$ then one calls $\widehat{u}(t)$ asymptotically stable. The most important element of the proof of this analysis is the study of the spectral measure of the operator one obtains by linearizing around $\alpha_{\mu}$. One then considers the associated linearized NLS.

If $u=e^{-i \theta}\left(\alpha_{\mu_{1}}+\beta\right)$ satisfies the NLS then $\beta$ satisfies the linearized NLS (LNLS)

$$
\begin{equation*}
i \mathrm{~d}_{t} \vec{\beta}=H \vec{\beta}+\vec{\gamma}, \tag{298}
\end{equation*}
$$

where

$$
\begin{align*}
& H:=\left(L_{0}+\mu\right) D-(\sigma+1) \alpha^{2 \sigma} D-\sigma \alpha^{2 \sigma} J,  \tag{299}\\
& D:=\left[\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right], \quad J:=\left[\begin{array}{cc}
0 & 1 \\
-1 & 0
\end{array}\right], \quad \vec{\beta}:=\left[\begin{array}{c}
\beta \\
\bar{\beta}
\end{array}\right], \quad \vec{\gamma}:=\left[\begin{array}{c}
\gamma \\
-\bar{\gamma}
\end{array}\right], \tag{300}
\end{align*}
$$

$$
\begin{align*}
\alpha_{\mu_{1}}=\alpha & +\widehat{\alpha}, \quad \alpha_{\mu}=\alpha, \quad \gamma=\gamma_{0}+\gamma_{1}, \quad \gamma_{0}:=-\mathrm{d}_{t} \nu_{1} \alpha_{\mu_{1}}+i \partial_{\mu_{1}} \alpha_{\mu_{1}} \mathrm{~d}_{t} \mu_{1}  \tag{301}\\
\gamma_{1}:=\widehat{\mu} \beta & -[(\sigma+1) \beta+\sigma \bar{\beta}] \sum_{j=1}^{2 \sigma}\binom{2 \sigma}{j} \widehat{\alpha}^{j} \alpha^{2 \sigma-j}-\mathrm{d}_{t} \nu_{1} \beta  \tag{302}\\
& +\sum_{(j, k)}^{\prime}\binom{\sigma}{j} \bar{\beta}^{j} \alpha_{\mu_{1}}^{\sigma-j}\binom{\sigma+1}{k} \beta^{k} \alpha_{\mu_{1}}^{\sigma+1-k}, \tag{303}
\end{align*}
$$

where $\sum_{(j, k)}^{\prime}$ sums over all $(j, k) \in \mathbb{Z} \times \mathbb{Z}$ for $0 \leq j \leq \sigma$ and $0 \leq k \leq \sigma+1$ with the exclusion of $(0,0),(0,1)$ and $(1,0)$.

We will arrive at the properties of $H$ by studying a sequence of simpler operators.

## Definition 1.1.

$$
\begin{align*}
& H_{0}:=\left(L_{0}+\mu\right) D, \quad H_{1}:=H_{0}-q_{1} \Pi_{0} D, \quad H_{2}:=H_{1}-q_{2} \Pi_{0} J,  \tag{304}\\
& \rho:=\alpha(0), \quad q_{1}:=(\sigma+1) \rho^{2 \sigma}, \quad q_{2}:=\sigma \rho^{2 \sigma},  \tag{305}\\
& L:=L_{0}-q_{1} \Pi_{0}, \quad U:=H-H_{2} . \tag{306}
\end{align*}
$$

Although $H_{0}$ is self-adjoint, it is the case that $H_{2}$ and $H$ are not. This property of $H$ is typical of linearized operators and makes the analysis very difficult for most systems.

Theorem 5. The spectrum of $H_{2}$ has the following properties.
(1) $\sigma_{\mathrm{d}}\left(H_{2}\right)=\sigma_{\mathrm{p}}\left(H_{2}\right)=\left\{(-1)^{j} i(2 \sigma)^{1 / 2} \mu^{-\sigma}\left[1+\mathcal{O}\left(\mu^{-1}\right)\right]\right\}_{j=0}^{1}$.
(2) $\sigma_{\mathrm{e}}\left(H_{2}\right)=\sigma_{\mathrm{c}}\left(H_{2}\right)=\sigma_{\mathrm{ac}}\left(H_{2}\right)=(-\infty,-\mu] \cup[\mu, \infty)$.

Theorem 6. The spectrum of $H$ has the following properties.
(1) $\sigma_{\mathrm{d}}(H)=\sigma_{\mathrm{p}}(H)=\{0\}$, with multiplicity 2 .
(2) $\sigma_{\mathrm{e}}(H)=\sigma_{\mathrm{c}}(H)=\sigma_{\mathrm{ac}}(H)=(-\infty,-\mu] \cup[\mu, \infty)$.

Theorem 7. For all $-3 \geq \tau \in \mathbb{R}, v \in \ell^{1}$, there exists a constant $1<\kappa \in \mathbb{R}$ such that

$$
\begin{equation*}
\left\|W_{\kappa, \tau} e^{-i t H_{2}} \Pi_{\mathrm{e}}^{H_{2}} W_{\kappa, \tau} v\right\|_{\infty}=\mathcal{O}\left(t^{-1} \log ^{-2} t\right), \quad t \nearrow \infty \tag{307}
\end{equation*}
$$

Theorem 8. For all $-3 \geq \tau \in \mathbb{R}, W_{\kappa, \tau}^{-1} v \in \ell^{1}$, there exists a constant $1<\kappa \in \mathbb{R}$ such that

$$
\begin{equation*}
\left\|W_{\kappa, \tau} e^{-i t H} \Pi_{\mathrm{e}}^{H} W_{\kappa, \tau} v\right\|_{\infty}=\mathcal{O}\left(t^{-1} \log ^{-2} t\right), \quad t \nearrow \infty \tag{308}
\end{equation*}
$$

One may extract from $E_{1}(z):=\int_{1}^{\infty} \mathrm{d} t e^{-z t} t^{-1}$ the well-known asymptotic expansion

$$
\begin{equation*}
E_{1}(z):=e^{-z} z^{-1} \sum_{k=1}^{n-1} \frac{k!}{(-z)^{k}}+\mathcal{O}\left(n!z^{-n}\right) \tag{309}
\end{equation*}
$$

which is valid for large values of $\Re z$. This in turn gives

$$
\begin{equation*}
\psi_{z}(0):=z^{-1} \sum_{k=1}^{n-1} \frac{k!}{(-z)^{k}}+\mathcal{O}\left(n!z^{-n}\right) \tag{310}
\end{equation*}
$$

which will be a crucial tool in our analysis.

## 2. Spectral Properties of $\mathrm{H}_{2}$

Consider that $A$ is an essentially self-adjoint operator on $\mathscr{H}$ and $B \in \mathcal{B}(\mathscr{H})$. If $A$ and $B$ commute strongly on their common domain then $R_{z}^{A}$ commutes strongly with $B$ on $\mathscr{H}$ for all $z \in \rho(A)$. Furthermore consider $H=A D$ where $D$ is the diagonal matrix defined above, it is the cas that $R_{z}^{H}=\left[\begin{array}{cc}R_{z}^{A} & 0 \\ 0 & -R_{-z}^{A}\end{array}\right]$ since

$$
\begin{align*}
H-z & =\left[\begin{array}{cc}
A-z & 0 \\
0 & -A-z
\end{array}\right]=\left[\begin{array}{cc}
A-z & 0 \\
0 & -(A+z)
\end{array}\right]  \tag{311}\\
\Rightarrow(H-z)^{-1} & =\left[\begin{array}{cc}
(A-z)^{-1} & 0 \\
0 & -(A+z)^{-1}
\end{array}\right] \tag{312}
\end{align*}
$$

One may write $R_{z}^{H_{1}}=\left[\begin{array}{cc}R_{z_{1}}^{L} & 0 \\ 0 & -R_{z_{2}}^{L}\end{array}\right]$, where here $L=L_{0}-q_{1} \Pi_{0}$ and where $q_{1}, z_{1}, z_{2}$ are defined as given above. Since $H_{2}^{2}$ and $H^{2}$ are self-adjoint, it follows that $\sigma\left(H_{2}\right), \sigma(H) \subseteq$ $\mathbb{R} \cup i \mathbb{R}$, see e.g. [25].

The first part of our analysis will be dedicated to proving that the point spectrum of $H_{2}$ consists of a conjugate pair of complex eigenvalues which are very close to the origin for large $\mu$.

Lemma 2.1. The eigenvalues of $H_{2}$ are given by the roots of $h(z):=q_{2}^{2} f_{z_{1}}^{L} f_{z_{2}}^{L}-1$, where $z_{1}:=z-\mu$ and $z_{2}:=-z-\mu$.

Proof. Let $\vec{v}=\left[\begin{array}{l}v_{1} \\ v_{2}\end{array}\right]$ and take $v_{1}(0) \neq 0$.

$$
\begin{align*}
H_{2} \vec{v} & =z \vec{v}  \tag{313}\\
\left(H_{1}-q_{2} \Pi_{0} J\right) \vec{v} & =  \tag{314}\\
\Rightarrow \quad \vec{v} & =R_{z}^{H_{1}} q_{2} \Pi_{0} J \vec{v}  \tag{315}\\
\Rightarrow \quad\left(\chi_{0}, \vec{v}\right) & =\left(\chi_{0}, R_{z}^{H_{1}} q_{2} \Pi_{0} J \vec{v}\right) \tag{316}
\end{align*}
$$

$$
\begin{align*}
\vec{v}(0) & =q_{2} f_{z}^{H_{1}} J \vec{v}(0)  \tag{317}\\
{\left[\begin{array}{l}
v_{1}(0) \\
v_{2}(0)
\end{array}\right] } & =q_{2}\left[\begin{array}{cc}
0 & f_{z_{1}}^{L} \\
f_{z_{2}}^{L} & 0
\end{array}\right]\left[\begin{array}{l}
v_{1}(0) \\
v_{2}(0)
\end{array}\right]  \tag{318}\\
\Rightarrow \quad v_{1}(0) & =q_{2}^{2} f_{z_{1}}^{L} f_{z_{2}}^{L} v_{1}(0)  \tag{319}\\
\Rightarrow \quad 1 & =q_{2}^{2} f_{z_{1}}^{L} f_{z_{2}}^{L} . \tag{320}
\end{align*}
$$

Next we find some preliminary estimates which are asymptotic in $\mu \nearrow \infty$.

Lemma 2.2. Let $\widehat{\epsilon}:=\left(\psi_{-\mu},\left(I-\Pi_{0}\right) \alpha^{2 \sigma+1}\right)$. It is the case that $\rho^{-1} \widehat{\epsilon}=\mu^{-(2 \sigma+2)}+$ $\mathcal{O}\left(a^{-(2 \sigma+3)}\right)$ for all $\sigma \in \mathbb{Z}_{+}$.

Proof. One has that $\left(L_{0}+\mu\right) \alpha=\alpha^{2 \sigma+1} \Rightarrow \alpha=R_{-\mu}^{L_{0}} \alpha^{2 \sigma+1} \Rightarrow 1=\rho^{2 \sigma} f_{-\mu}+\rho^{-1} \widehat{\epsilon}$.
Let $\epsilon_{x}:=\alpha(x) / \alpha(x-1)$, for all $0<x \in \mathbb{Z}$. One may observe that

$$
\begin{align*}
&-\mu \alpha+\alpha^{2 \sigma+1}= L_{0} \alpha  \tag{321}\\
& \Rightarrow \quad-\mu \alpha(x)+\alpha^{2 \sigma+1}(x)=-(x+1) \alpha(x+1)  \tag{322}\\
&+(2 x+1) \alpha(x)-x \alpha(x-1)  \tag{323}\\
& \Rightarrow \quad x \epsilon_{1} \ldots \epsilon_{x-1}+\left(\rho \epsilon_{1} \ldots \epsilon_{x}\right)^{2 \sigma+1}=-(x+1) \rho \epsilon_{1} \ldots \epsilon_{x+1}  \tag{324}\\
&+(2 x+1+\mu) \rho \epsilon_{1} \ldots \epsilon_{x}  \tag{325}\\
& \Rightarrow \quad x+\left(\rho \epsilon_{1} \ldots \epsilon_{x-1}\right)^{2 \sigma} \epsilon_{x}^{2 \sigma+1}=-(x+1) \epsilon_{x} \epsilon_{x+1}+(2 x+1+\mu) \epsilon_{x}  \tag{326}\\
& \Rightarrow \quad \epsilon_{x}=\left[\mu+1+2 x-(x+1) \epsilon_{x+1}\right]^{-1}\left[x+\left(\rho \epsilon_{1} \ldots \epsilon_{x-1}\right)^{2 \sigma} \epsilon_{x}^{2 \sigma+1}\right] \tag{327}
\end{align*}
$$

The analogous equation for $x=0$ is

$$
\begin{equation*}
\rho^{2 \sigma}=\mu+1-\epsilon_{1} . \tag{328}
\end{equation*}
$$

We recall that $\sum_{x=0}^{\infty} \psi_{-\mu}(x)=\mu^{-1}$. Therefore

$$
\begin{align*}
\rho^{-1} \widehat{\epsilon} & =\rho^{-1} \sum_{x=1}^{\infty} \psi_{-\mu}(x) \alpha^{2 \sigma+1}(x)<\rho^{-1} \sum_{x=1}^{\infty} \psi_{-\mu}(x) \alpha^{2 \sigma+1}(1)  \tag{331}\\
& =\rho^{2 \sigma} \epsilon_{1}^{2 \sigma+1} \sum_{x=1}^{\infty} \psi_{-\mu}(x)<(\mu+1)\left[\mu^{-1}+\mathcal{O}\left(\mu^{-2}\right)\right]^{2 \sigma+1}\left[\mu^{-2}+\mathcal{O}\left(\mu^{-3}\right)\right]  \tag{332}\\
& <\mu^{-(2 \sigma+2)}+\mathcal{O}\left(\mu^{-(2 \sigma+3)}\right) \tag{333}
\end{align*}
$$

Since the lowest order of this bound has a power which depends explicitly on $\sigma$ and should hold for all $0<\sigma \in \mathbb{Z}$ it follows that the full asymptotic expansion, as well as the error term, must consist of orders which are powers of $\mu^{-1}$ that depend explicitly on $\sigma$.

$$
\begin{align*}
\psi_{-\mu}(1)= & (1+\mu) \psi_{-\mu}(0)-1  \tag{334}\\
= & (1+\mu)\left[\mu^{-1}-\mu^{-2}+2 \mu^{-3}+\mathcal{O}\left(\mu^{-4}\right)\right]-1=\mu^{-2}+\mathcal{O}\left(\mu^{-3}\right),  \tag{335}\\
\rho^{-1} \widehat{\epsilon}= & \rho^{-1} \sum_{x=1}^{\infty} \psi_{-\mu}(x) \alpha^{2 \sigma+1}(x)  \tag{336}\\
= & \psi_{-\mu}(1) \rho^{2 \sigma} \epsilon_{1}^{2 \sigma+1}+\rho^{-1} \sum_{x=2}^{\infty} \psi_{-\mu}(x) \alpha^{2 \sigma+1}(x)  \tag{337}\\
= & {\left[\mu^{-2}+\mathcal{O}\left(\mu^{-3}\right)\right]\left[\mu+1+\mathcal{O}\left(\mu^{-1}\right)\right]\left[\mu^{-1}+\mathcal{O}\left(\mu^{-2}\right)\right]^{2 \sigma+1} }  \tag{338}\\
& \quad+\rho^{-1} \sum_{x=2}^{\infty} \psi_{-\mu}(x) \alpha^{2 \sigma+1}(x)  \tag{339}\\
= & \mu^{-(2 \sigma+2)}+\mathcal{O}\left(\mu^{-(2 \sigma+3)}\right)+\rho^{-1} \sum_{x=2}^{\infty} \psi_{-\mu}(x) \alpha^{2 \sigma+1}(x) \tag{340}
\end{align*}
$$

$$
\begin{align*}
& \Rightarrow \quad \rho^{-1} \sum_{x=2}^{\infty} \psi_{-\mu}(x) \alpha^{2 \sigma+1}(x)=\mathcal{O}\left(\mu^{-(2 \sigma+3)}\right)  \tag{341}\\
& \Rightarrow \quad \rho^{-1} \widehat{\epsilon}=\mu^{-(2 \sigma+2)}+\mathcal{O}\left(\mu^{-(2 \sigma+3)}\right) \tag{342}
\end{align*}
$$

Now we show that there are no real eigenvalues through a series of lemmas.

Lemma 2.3. It is the case that $h(0)=2 \sigma^{-1} \mu^{-(2 \sigma+2)}+\mathcal{O}\left(\mu^{-(2 \sigma+3)}\right)>0$.

Proof. One may observe that

$$
\begin{gather*}
\rho^{2 \sigma} f_{-\mu}=1-\mu^{-(2 \sigma+2)}+\mathcal{O}\left(\mu^{-(2 \sigma+3)}\right)  \tag{343}\\
f_{-\mu}^{-1}=\mu+1-\mu^{-1}+3 \mu^{-2}-16 \mu^{-3}+\mathcal{O}\left(\mu^{-4}\right)  \tag{344}\\
\rho^{2 \sigma}=\mu+1-\mu^{-1}+3 \mu^{-2}-16 \mu^{-3}+\mathcal{O}\left(\mu^{-4}\right)  \tag{345}\\
\quad-\mu^{-(2 \sigma+1)}-\mu^{-(2 \sigma+2)}+\mathcal{O}\left(\mu^{-(2 \sigma+3)}\right)  \tag{346}\\
\rho^{2 \sigma} f_{-2 \mu}=2^{-1}\left[1+2^{-1} \mu^{-1}-\mu^{-2}+\mathcal{O}\left(\mu^{-3}\right)\right]  \tag{347}\\
q_{1} f_{-2 \mu}-1=2^{-1}(\sigma-1)+4^{-1}(\sigma+1) \mu^{-1}-4^{-1}(\sigma+1) \mu^{-2}+\mathcal{O}\left(\mu^{-3}\right)  \tag{348}\\
h(0)=\left[\left(q_{1} f_{-\mu}-1\right)^{-1} q_{2} f_{-\mu}\right]^{2}-1 \\
=\left\{\sigma^{-1}\left[1-\left(1+\sigma^{-1}\right) \mu^{-(2 \sigma+2)}+\mathcal{O}\left(\mu^{-(2 \sigma+3)}\right)\right]^{-1} \sigma \rho^{2 \sigma} f_{-\mu}\right\}^{2}-1 \\
=\left\{\left[1+\left(1+\sigma^{-1}\right) \mu^{-(2 \sigma+2)}+\mathcal{O}\left(\mu^{-(2 \sigma+3)}\right)\right]\right. \\
\left.\quad \times\left[1-\mu^{-(2 \sigma+2)}+\mathcal{O}\left(\mu^{-(2 \sigma+3)}\right)\right]\right\}^{2}-1 \\
= \\
\\
=\left[1+\sigma^{-1} \mu^{-(2 \sigma+2)}+\mathcal{O}\left(\mu^{-(2 \sigma+3)}\right)\right]^{2}-1 \\
=2 \sigma^{-1} \mu^{-(2 \sigma+2)}+\mathcal{O}\left(\mu^{-(2 \sigma+3)}\right)>0 .
\end{gather*}
$$

Lemma 2.4. For $z=a \in[-\mu, \mu]$, we define $a_{1}:=a-\mu, a_{2}:=-a-\mu$. One has for $\sigma=1$

$$
\begin{equation*}
\lim _{a \nearrow \mu} h(a)=2^{-1} \mu-4^{-1}+\mathcal{O}\left(\mu^{-1}\right)>0 \tag{355}
\end{equation*}
$$

and for $1<\sigma \in \mathbb{Z}_{+}$

$$
\begin{align*}
\lim _{a \nearrow \mu} h(a)=( & \left.\sigma^{2}-1\right)^{-1}+(\sigma+1)^{-1}(\sigma-1)^{-2} \sigma^{2} \mu^{-1}  \tag{356}\\
& +4^{-1}(\sigma+1)^{-1}(\sigma-1)^{-3}\left(3 \sigma^{2}+10 \sigma-9\right) \mu^{-2}+\mathcal{O}\left(\mu^{-3}\right)>0 \tag{357}
\end{align*}
$$

Proof. Let $s_{1}:=-a_{1} / a, s_{1} \in[0,2]$. For $a_{1} \gg 1$ one has

$$
\begin{align*}
(\sigma+ & 1)^{-1} s_{1}\left(q_{1} f_{a_{1}}-1\right)=\left[1-(\sigma+1)^{-1} s_{1}\right]+\left(1-s_{1}^{-1}\right) \mu^{-1}  \tag{358}\\
& -\left(1+s_{1}^{-1}-2 s_{1}^{-2}\right) \mu^{-2}+\left(16+3 s_{1}^{-1}+2 s_{1}^{-2}+6 s_{1}^{-3}-24 s_{1}^{-4}\right) \mu^{-4}  \tag{359}\\
& +\mathcal{O}\left(\mu^{-5}\right)-\mu^{-(2 \sigma+2)}-\left(1-s_{1}^{-1}\right) \mu^{-(2 \sigma+3)}+\mathcal{O}\left(\mu^{-(2 \sigma+4)}\right) \tag{360}
\end{align*}
$$

We recall that the generalized exponential integrals have the representation

$$
\begin{align*}
E_{n+1}(z)= & -\frac{(-z)^{n}}{n!} \log z+\frac{e^{-z}}{n!} \sum_{k=1}^{n}(-z)^{k-1}(n-k)!  \tag{361}\\
& +\frac{e^{-z}(-z)^{n}}{n!} \sum_{k=0}^{\infty} \frac{z^{k}}{k!} \digamma(k+1) \tag{362}
\end{align*}
$$

This permits one to write

$$
\begin{equation*}
f_{-a}=e^{a} E_{1}(a)=-e^{a} \log a+\sum_{k=0}^{\infty} \frac{a^{k}}{k!} \digamma(k+1) \tag{363}
\end{equation*}
$$

and hereby one may observe that near $a=0$ the above expression for $f_{-a}$ is dominated by logarithmic behavior. Therefore $\lim _{a \not \nearrow_{\mu}}\left(-f_{a_{1}}^{L}\right)=\lim _{a \not ~_{\mu}}\left(q_{1} f_{a_{1}}-1\right)^{-1} f_{a_{1}}=q_{1}^{-1}$.
(364) $\quad \lim _{a \nearrow \mu} h(a)=\lim _{a \nearrow \mu}\left(q_{1} f_{z_{1}}-1\right)^{-1} q_{2} f_{z_{1}}\left(q_{1} f_{z_{2}}-1\right)^{-1} q_{2} f_{z_{2}}-1$

$$
\begin{align*}
& =q_{1}^{-1} q_{2}^{2}\left(q_{1} f_{-2 \mu}-1\right)^{-1} f_{-2 \mu}-1  \tag{365}\\
& =(\sigma+1)^{-1} \sigma^{2}\left[(\sigma+1) \rho^{2 \sigma} f_{-2 \mu}-1\right]^{-1} \rho^{2 \sigma} f_{-2 \mu}-1  \tag{366}\\
& =2^{-1}(\sigma+1)^{-1} \sigma^{2}\left[2^{-1}(\sigma-1)+4^{-1}(\sigma+1) \mu^{-1}\right. \tag{367}
\end{align*}
$$

$$
\begin{equation*}
\left.-4^{-1}(\sigma+1) \mu^{-2}+\mathcal{O}\left(\mu^{-3}\right)\right]^{-1}\left[1+2^{-1} \mu^{-1}-\mu^{-2}+\mathcal{O}\left(\mu^{-3}\right)\right] \tag{368}
\end{equation*}
$$

For $\sigma=1$ :
$\lim _{a \nearrow \mu} h(a)=4^{-1}\left[2^{-1} \mu^{-1}-2^{-1} \mu^{-2}+\mathcal{O}\left(\mu^{-3}\right)\right]^{-1}\left[1+2^{-1} \mu^{-1}-\mu^{-2}+\mathcal{O}\left(\mu^{-3}\right)\right]-1$

$$
\begin{align*}
& =2^{-1} \mu\left[1-\mu^{-1}+\mathcal{O}\left(\mu^{-2}\right)\right]^{-1}\left[1+2^{-1} \mu^{-1}-\mu^{-2}+\mathcal{O}\left(\mu^{-3}\right)\right]-1  \tag{370}\\
& =2^{-1} \mu\left[1+\mu^{-1}+\mathcal{O}\left(\mu^{-2}\right)\right]\left[1+2^{-1} \mu^{-1}-\mu^{-2}+\mathcal{O}\left(\mu^{-3}\right)\right]-1  \tag{371}\\
& =2^{-1} \mu-4^{-1}+\mathcal{O}\left(\mu^{-1}\right)>0 \tag{372}
\end{align*}
$$

For $\sigma>1$ :

$$
\begin{align*}
& \lim _{a \nearrow \mu} h(a)=( \left.\sigma^{2}-1\right)^{-1} \sigma^{2}\left\{1-2^{-1}(\sigma-1)^{-1}(\sigma+1)\left[-\mu^{-1}+\mu^{-2}+\mathcal{O}\left(\mu^{-3}\right)\right]\right\}^{-1}  \tag{373}\\
& \times\left[1+2^{-1} \mu^{-1}-\mu^{-2}+\mathcal{O}\left(\mu^{-3}\right)\right]-1  \tag{374}\\
&=( \left.\sigma^{2}-1\right)^{-1} \sigma^{2}\left\{1+2^{-1}(\sigma-1)^{-1}(\sigma+1)\left[-\mu^{-1}+\mu^{-2}+\mathcal{O}\left(\mu^{-3}\right)\right.\right.  \tag{375}\\
&\left.+\left[2^{-1}(\sigma-1)^{-1}(\sigma+1)\right]^{2}\left[-\mu^{-1}+\mathcal{O}\left(\mu^{-2}\right)\right]^{2}+\mathcal{O}\left(\mu^{-3}\right)\right\}  \tag{376}\\
& \times\left[1+2^{-1} \mu^{-1}-\mu^{-2}+\mathcal{O}\left(\mu^{-3}\right)\right]-1  \tag{377}\\
&=\left(\sigma^{2}-1\right)^{-1} \sigma^{2}\left[1-2^{-1}(\sigma-1)^{-1}(\sigma+1) \mu^{-1}\right.  \tag{378}\\
&\left.+4^{-1}(\sigma-1)^{2}(\sigma+1)(3 \sigma-1) \mu^{-2}+\mathcal{O}\left(\mu^{-3}\right)\right]  \tag{379}\\
& \times\left[1+2^{-1} \mu^{-1}-\mu^{-2}+\mathcal{O}\left(\mu^{-3}\right)\right]-1  \tag{380}\\
&=( \left.\sigma^{2}-1\right)^{-1} \sigma^{2}\left[1-(\sigma-1)^{-1} \mu^{-1}\right.  \tag{381}\\
&\left.+4^{-1}(\sigma-1)^{2}\left(3 \sigma^{2} 10 \sigma-9\right) \mu^{-2}+\mathcal{O}\left(\mu^{-3}\right)\right]-1  \tag{382}\\
&=( \left.\sigma^{2}-1\right)^{-1}-(\sigma+1)^{-1}(\sigma-1)^{-2} \sigma^{2} \mu^{-1}  \tag{383}\\
&+4^{-1}(\sigma+1)^{-1}(\sigma-1)^{3}\left(3 \sigma^{2}+10 \sigma-9\right) \mu^{-2}+\mathcal{O}\left(\mu^{-3}\right)>0 \tag{384}
\end{align*}
$$

Lemma 2.5. Let $h_{0}(z):=\left(q_{1} f_{z_{1}}-1\right)^{-1}\left(q_{1} f_{z_{2}}-1\right)^{-1}(2 \sigma+1) \rho^{4 \sigma}$. It is the case that $h_{0}(z)>0$ for all $z=a \in[-\mu, \mu]$.

Proof. Let $c_{0}:=(2 \sigma+1)^{-1} \rho^{-4 \sigma}, c_{1}:=(2 \sigma+1)^{-1}(\sigma+1) \rho^{-2 \sigma}, c_{2}:=c_{1}^{2}-c_{0}=$ $(2 \sigma+1)^{-2} \rho^{-4 \sigma} \sigma^{2}$. One then has $h(z)=h_{0}(z)\left[c_{2}-\left(f_{z_{1}}-c_{1}\right)\left(f_{z_{2}}-c_{1}\right)\right]$. For all $\sigma \in \mathbb{Z}_{+}$and $a \in[-\mu, \mu]$ it is the case that $q_{1} f_{a_{1}}=q_{1} f_{a-\mu}$ and $q_{1} f_{a_{2}}=q_{1} f_{-a-\mu}$ vary monotonically between

$$
\begin{equation*}
q_{1} f_{-2 \mu}=2^{-1}(\sigma+1)+4^{-1}(\sigma+1) \mu^{-1}-4^{-1}(\sigma+1) \mu^{-2}+\mathcal{O}\left(\mu^{-3}\right)>1 \tag{385}
\end{equation*}
$$

and

$$
\begin{equation*}
\lim _{a \nearrow 0} q_{1} f_{-a}=\infty>1 \tag{386}
\end{equation*}
$$

Therefore $\left(q_{1} f_{a_{i}}-1\right)^{-1}>0$ which concludes the proof.

Lemma 2.6. For $1 \ll(2 \sigma+1)^{-1} \mu=\mathcal{O}(\mu)$ it is the case that $\left(f_{a_{j}}-c_{1}\right), j=1,2$, have unique roots, $a=r_{j}:=(-1)^{j}(\sigma+1)^{-1}(\mu+\sigma)+\mathcal{O}\left(\mu^{-1}\right)$.

Proof. Consider $c_{1}=f_{a_{j}}$. Let $a_{1}=-s_{1} \mu, s_{1} \in[0,2]$, and assume that $1 \ll$ $s_{1} \mu=\mathcal{O}(\mu)$ so that the asymptotic expansion of $f_{a_{1}}$ is valid.

$$
\begin{align*}
c_{1}= & f_{a_{1}}=\left(s_{1} \mu\right)^{-1}-\left(s_{1} \mu\right)^{-2}+\mathcal{O}\left(\mu^{-3}\right)  \tag{387}\\
= & \left(s_{1} \mu\right)^{-1}\left[1-\left(s_{1} \mu\right)^{-1}+\mathcal{O}\left(\mu^{-2}\right)\right]  \tag{388}\\
\Rightarrow \quad s_{1}= & \mu^{-1} c_{1}^{-1}\left[1-\left(s_{1} \mu\right)^{-1}+\mathcal{O}\left(\mu^{-2}\right)\right]  \tag{389}\\
= & (\sigma+1)^{-1}(2 \sigma+1)\left[1+\mu^{-1}+\mathcal{O}\left(\mu^{-2}\right)\right]  \tag{390}\\
& \times\left[1-\left(s_{1} \mu\right)^{-1}+\mathcal{O}\left(\mu^{-2}\right)\right]  \tag{391}\\
= & (\sigma+1)^{-1}(2 \sigma+1)\left[1+\left(1-s^{-1}\right) \mu^{-1}+\mathcal{O}\left(\mu^{-2}\right)\right]  \tag{392}\\
\Rightarrow \quad s_{1}= & (\sigma+1)^{-1}(2 \sigma+1)  \tag{393}\\
& \times\left\{1+\left[1-(2 \sigma+1)^{-1}(\sigma+1)\right] \mu^{-1}+\mathcal{O}\left(\mu^{-2}\right)\right\}  \tag{394}\\
= & (\sigma+1)^{-1}(2 \sigma+1)+(\sigma+1)^{-1} \sigma \mu^{-1}+\mathcal{O}\left(\mu^{-2}\right) \tag{395}
\end{align*}
$$

This result satisfies the assumptions and the root must be unique, therefore the above value of $s_{1}$ specifies the unique root. In terms of $a$ one has

$$
\begin{align*}
a_{1} & =a-\mu=-s_{1} \mu  \tag{396}\\
\Rightarrow \quad a & =\left(1-s_{1}\right) \mu=\left\{1-\left[(\sigma+1)^{-1}(2 \sigma+1)+(\sigma+1)^{-1} \sigma \mu^{-1}+\mathcal{O}\left(\mu^{-2}\right)\right]\right\} \mu  \tag{397}\\
& =-(\sigma+1)^{-1}(\mu+\sigma)+\mathcal{O}\left(\mu^{-1}\right) \tag{398}
\end{align*}
$$

Due to symmetry between the $a_{j}$ there are two roots $a=r_{j}=(-1)^{j}(\sigma+1)^{-1}(\mu+$ $\sigma)+\mathcal{O}\left(\mu^{-1}\right)$.

Lemma 2.7. For $1 \ll r_{j}=\mathcal{O}(\mu)$ one has that $h(z)>0$ for all $z=a \in[-\mu, \mu]$.
Proof. Let $h_{1}(z):=\left(f_{z_{1}}-c_{1}\right)\left(f_{z_{2}}-c_{1}\right)$. One can observe that $h_{1}(0)=\left(f_{-\mu}-\right.$ $\left.c_{1}\right)^{2}>0$ and $\lim _{a \not \nearrow_{\mu}}\left(f_{a_{1}}-c_{1}\right)=\infty$. Then since $\left(f_{a_{j}}-c_{1}\right)$ are monotonic in $a$ and have unique roots $r_{i}$, it is the case that $h(a)>0$ for $a \in\left(r_{1}, r_{2}\right)$ and $h(a)<0$ for $a \in[-\mu, \mu] \backslash\left(r_{1}, r_{2}\right)$. We have shown that $h(0)>0$ and $h(\mu)=h(-\mu)>0$. Since $h(a)>0$ for $a \in\left(r_{1}, r_{2}\right)$ and $h(a)<0$ for $a \in[-\mu, \mu] \backslash\left(r_{1}, r_{2}\right)$ it must be the case that $h(a)>0$ for $a \in[-\mu, \mu] \backslash\left(r_{1}, r_{2}\right)$. It remains to consider $a \in\left(r_{1}, r_{2}\right)$. By assumption $1 \ll a=\mathcal{O}(\mu)$ and therefore the asymptotic expansion of the $f_{a_{j}}$ is valid.

$$
\begin{align*}
h_{1}(a)= & \left(f_{a_{1}}-c_{1}\right)\left(f_{a_{2}}-c_{1}\right)=f_{a_{1}} f_{a_{2}}-c_{1}\left(f_{a_{1}}+f_{a_{2}}\right)+c_{1}^{2}  \tag{399}\\
= & {\left[(\mu-a)^{-1}+\mathcal{O}\left(\mu^{-2}\right)\right]\left[(\mu+a)^{-1}+\mathcal{O}\left(\mu^{-2}\right)\right] }  \tag{400}\\
& -c_{1}\left\{\left[(\mu-a)^{-1}+\mathcal{O}\left(\mu^{-2}\right)\right]+\left[(\mu+a)^{-1}+\mathcal{O}\left(\mu^{-2}\right)\right]\right\}+c_{1}^{2}  \tag{401}\\
= & \left(\mu^{2}-a^{2}\right)^{-1}-\left(\mu^{2}-a^{2}\right)^{-1}(2 \sigma+1)^{-1}(\sigma+1)  \tag{402}\\
& \times\left[1+\mathcal{O}\left(\mu^{-1}\right)\right]\left[2+\mathcal{O}\left(\mu^{-1}\right)\right]+c_{1}^{2}+\mathcal{O}\left(\mu^{-3}\right)  \tag{403}\\
& =-\left(\mu^{2}-a^{2}\right)^{-1}(2 \sigma+1)^{-1}+c_{1}^{2}+\mathcal{O}\left(\mu^{-3}\right) \tag{404}
\end{align*}
$$

Therefore $-h(a)$ decays monotonically as $|a| \nearrow r_{2}$ for sufficiently large $\mu$. This guarantees that $h(a)>0$ for all $a \in\left(r_{1}, r_{2}\right)$.

Lemma 2.8. It is the case that $h(z)$ has no roots for $z \in(-\infty,-\mu] \cup[\mu, \infty)$.

Proof. This follows by the same principle which permits $L$ from having embedded eigenvalues.

Now we prove the existence and location (asymptotically) of the imaginary roots.

Lemma 2.9. $h(z)$ has exactly two roots, $\lambda_{ \pm}=z_{ \pm}:= \pm i(2 \sigma)^{1 / 2} \mu^{-\sigma}\left[1+\mathcal{O}\left(\mu^{-1}\right)\right]$, for $z \in i \mathbb{R}$.

Proof. One may observe that

$$
\begin{align*}
f_{z} & =(-z)^{-1}-(-z)^{-2}+2(-z)^{-3}-6(-z)^{-4}+\mathcal{O}\left(z^{-5}\right)  \tag{405}\\
f_{z}^{L} & =\left(1-q_{1} f_{z}\right)^{-1} f_{z}, \quad \partial_{z} f_{z}^{L}=\left(1-q_{1} f_{z}\right)^{-2} \partial_{z} f_{z}  \tag{406}\\
\partial_{z} f_{z} & =(-1)\left(f_{z}+z^{-1}\right), \quad \partial_{z}^{2} f_{z}=2 q_{1}\left(1-q_{1} f_{z}\right)^{-3}\left(\partial_{z} f_{z}\right)^{2}+\left(1-q_{1} f_{z}\right)^{-2} \partial_{z}^{2} f_{z},  \tag{407}\\
f_{-\mu}^{\prime} & =\mu^{-2}-2 \mu^{-3}+6 \mu^{-4}+\mathcal{O}\left(\mu^{-5}\right), \quad f_{-\mu}^{\prime \prime}=2 \mu^{-3}+6 \mu^{-4}+\mathcal{O}\left(\mu^{-5}\right) \tag{408}
\end{align*}
$$

and that for $z \in i \mathbb{R}$ one has that $h(z)=\left|q_{2} f_{\mu-z}^{L}\right|^{2}-1 \in \mathbb{R}$. Furthermore, for all $z \in i \mathbb{R}$ the asymptotic expansion of $f_{z_{j}}$ is valid since $\Re z_{j}=a \gg 1$.

First consider $|z| \ll 1$. One finds

$$
\begin{align*}
& \partial_{z}^{2} h(z)= q_{2}^{2}\left(\partial_{z}^{2} f_{z_{1}}^{L} f_{z_{2}}^{L}+2 \partial_{z} f_{z_{1}}^{L} \partial_{z} f_{z_{2}}^{L}+f_{z_{1}}^{L} \partial_{z}^{2} f_{z_{2}}^{L}\right)  \tag{409}\\
&= q_{2}^{2}\left(\partial_{z_{1}}^{2} f_{z_{1}}^{L} f_{z_{2}}^{L}-2 \partial_{z_{1}} f_{z_{1}}^{L} \partial_{z_{2}} f_{z_{2}}^{L}+f_{z_{1}}^{L} \partial_{z_{2}}^{2} f_{z_{2}}^{L}\right)  \tag{410}\\
&= q_{2}^{2}\left\{\left[2 q_{1}\left(1-q_{1} f_{z_{1}}\right)^{-3}\left(\partial_{z_{1}} f_{z_{1}}\right)^{2}+\left(1-q_{1} f_{z_{1}}\right)^{-2} \partial_{z_{1}}^{2} f_{z_{1}}\right]\right\}\left(1-q_{1} f_{z_{2}}\right)^{-1} f_{z_{2}}  \tag{411}\\
&-2\left(1-q_{1} f_{z_{1}}\right)^{-2} \partial_{z_{1}} f_{z_{1}}\left(1-q_{1} f_{z_{2}}\right)^{-2} \partial_{z_{2}} f_{z_{2}}  \tag{412}\\
&\left.+\left(1-q_{1} f_{z_{1}}\right)^{-1} f_{z_{1}}\left[2 q_{1}\left(1-q_{1} f_{z_{2}}\right)^{-3}\left(\partial_{z_{2}} f_{z_{2}}\right)^{2}+\left(1-q_{1} f_{z_{2}}\right)^{-2} \partial_{z_{2}}^{2} f_{z_{2}}\right]\right\}  \tag{413}\\
& h^{\prime \prime}(0)=q_{2}^{2}\left\{\left[2 q_{1}\left(1-q_{1} f_{-\mu}\right)^{-3}\left(f_{-\mu}^{\prime}\right)^{2}+\left(1-q_{1} f_{-\mu}\right)^{-2} f_{-\mu}^{\prime \prime}\right]\right\}\left(1-q_{1} f_{-\mu}\right)^{-1} f_{-\mu}  \tag{414}\\
&-2\left(1-q_{1} f_{-\mu}\right)^{-2} f_{-\mu}^{\prime}\left(1-q_{1} f_{-\mu}\right)^{-2} f_{-\mu}^{\prime}  \tag{415}\\
&\left.+\left(1-q_{1} f_{-\mu}\right)^{-1} f_{-\mu}\left[2 q_{1}\left(1-q_{1} f_{-\mu}\right)^{-3}\left(f_{-\mu}^{\prime}\right)^{2}+\left(1-q_{1} f_{-\mu}\right)^{-2} f_{-\mu}^{\prime \prime}\right]\right\} \tag{416}
\end{align*}
$$

$$
\begin{align*}
&=q_{2}^{2} {\left[2\left(q_{1} f_{-\mu}-1\right)^{-4}\left(f_{-\mu}^{\prime}\right)^{2}\left(2 q_{1} f_{-\mu}-1\right)-2\left(q_{1} f_{-\mu}-1\right)^{-3} f_{-\mu}^{\prime \prime} f_{-\mu}\right] }  \tag{417}\\
&=q_{2}^{2}\left(2\left\{\sigma\left[1+\left(1+\sigma^{-1}\right) \mu^{-(2 \sigma+2)}+\mathcal{O}\left(\mu^{-(2 \sigma+3)}\right)\right]\right\}^{-4}\right.  \tag{418}\\
& \times\left[\mu^{-2}-2 \mu^{-3}+6 \mu^{-4}+\mathcal{O}\left(\mu^{-6}\right)\right]^{2}  \tag{419}\\
& \times\left\{2\left[(\sigma+1)-(\sigma+1) \mu^{-(2 \sigma+2)}+\mathcal{O}\left(\mu^{-(2 \sigma+3)}\right)\right]-1\right\}  \tag{420}\\
&-2\left\{\sigma\left[1+\left(1+\sigma^{-1}\right) \mu^{-(2 \sigma+2)}+\mathcal{O}\left(\mu^{-(2 \sigma+3)}\right)\right]\right\}^{-3}  \tag{421}\\
& \times\left[2 \mu^{-3}-6 \mu^{-4}+\mathcal{O}\left(\mu^{-5}\right)\right]  \tag{422}\\
&\left.\times\left[\mu^{-1}-\mu^{-2}+2 \mu^{-3}-6 \mu^{-4}+\mathcal{O}\left(\mu^{-5}\right)\right]\right)  \tag{423}\\
&=2 \sigma^{-2} \mu^{-2}\left[1-2 \mu^{-1}+\mathcal{O}\left(\mu^{-2}\right)\right] . \tag{424}
\end{align*}
$$

Assume that $h(z)=h(0)+2^{-1} h^{\prime \prime}(0) z^{2}+\epsilon$ where $|\epsilon| \leq \mathcal{O}\left(\mu^{-4}\right)$. One finds
(425) $h(z)=0$

$$
\begin{align*}
\Rightarrow z^{2} & =-2[h(0)+\epsilon]\left[h^{\prime \prime}(0)\right]^{-1}  \tag{426}\\
& =-(2)\left[2 \sigma^{-1} \mu^{-(2 \sigma+2)}+\mathcal{O}\left(\mu^{-(2 \sigma+3)}\right)+\epsilon\right]\left(2^{-1} \sigma^{2} \mu^{2}\right)\left[1-2 \mu^{-1}+\mathcal{O}\left(\mu^{-2}\right)\right]^{-1}  \tag{427}\\
& =-2 \sigma \mu^{2 \sigma}\left[1+\mathcal{O}\left(\mu^{-1}\right)+\epsilon\right] . \tag{428}
\end{align*}
$$

This result is compatible with the assumptions, thus there are at least the two imaginary roots given by $z_{ \pm}= \pm i(2 \sigma)^{1 / 2} \mu^{-\sigma}\left[1+\mathcal{O}\left(\mu^{-1}\right)\right]$.

It remains to be shown that there are no other imaginary roots. It is sufficient to prove that $h(z)$ has nonpositive curvature for all $z \in i \mathbb{R}$. Let $z_{j}=(-1)^{j+1} i b-\mu$, where $b \in \mathbb{R}$.

$$
\begin{align*}
h(z) & =q_{2}^{2}\left(q_{1}-f_{z_{1}}^{-1}\right)\left(q_{1}-f_{z_{2}}^{-1}\right)-1  \tag{429}\\
& =q_{2}^{2}[(\sigma+1) \mu-(\mu-i b)+\mathcal{O}(1)]^{-1}[(\sigma+1) \mu-(\mu+i b)+\mathcal{O}(1)]^{-1}-1  \tag{430}\\
& =q_{2}^{2}\left(\sigma^{2} \mu^{2}+b^{2}\right)^{-1}\left[1+\mathcal{O}\left(\mu^{-1}\right)\right]-1 \tag{431}
\end{align*}
$$

Therefore $h(i b)$ decays monotonically as $|b| \nearrow \infty$ and there are only two imaginary roots.

Proof of Theorem (5) Part (1). We have exhaustively shown that

$$
\begin{equation*}
\lambda_{ \pm}:= \pm i(2 \sigma)^{1 / 2} \mu^{-\sigma}\left[1+\mathcal{O}\left(\mu^{-1}\right)\right] \tag{432}
\end{equation*}
$$

are the only roots of $h(z)$ for $z \in \mathbb{R} \cup i \mathbb{R}$. The absence of embedded eigenvalues follows from arguments similar to those for $\sigma(L)$.

Proof of Theorem (5) Part (2). By Weyl's critereon it is the case that $\sigma_{\mathrm{e}}\left(H_{2}\right)=$ $\sigma_{\mathrm{e}}\left(H_{0}\right)=(-\infty,-\mu] \cup[\mu, \infty)$. It is clear that there exists a well-defined absolutely continuous spectral measure on $\sigma_{\mathrm{e}}\left(H_{2}\right)$. The representation of $H_{2} v=z v$ as a coupled series of algebraic equations guarantees that each $\lambda \in \sigma_{\mathrm{e}}\left(H_{2}\right)$ has multiplicity 1 . Therefore one must have that $\sigma_{\mathrm{e}}\left(H_{2}\right)=\sigma_{\mathrm{ac}}\left(H_{0}\right)$.

## 3. Spectral Properties of $H$

We recall Proposition (2) from Chapter 2:

$$
\begin{equation*}
\left\|\left(I-\Pi_{0}\right) \alpha_{\mu}\right\|_{1} \leq \mu^{-(2 \sigma)^{-1}(2 \sigma-1)}+\mathcal{O}\left(\mu^{-(2 \sigma)^{-1}(4 \sigma-1)}\right) \tag{433}
\end{equation*}
$$

This gives

$$
\begin{align*}
\|U\| & \leq 2(2 \sigma+1)\left\|\left(I-\Pi_{0}\right) \alpha^{2 \sigma}\right\|_{1} \leq 2(2 \sigma+1)\left\|\left(I-\Pi_{0}\right) \alpha\right\|_{1}  \tag{434}\\
& \leq 2(2 \sigma+1) \mu^{-(2 \sigma)^{-1}(2 \sigma-1)}+\mathcal{O}\left(\mu^{-(2 \sigma)^{-1}(4 \sigma-1)}\right)=: m(\mu) \tag{435}
\end{align*}
$$

We recall without proof a statement of Kato [16] regarding norm resolvent convergence.

Proposition. For $A$ a closed operator and $\left\{A_{n}\right\}_{n=0}^{\infty}$ a sequence of closed operators, if $R_{z}^{A_{n}}$ converges in norm to $R_{z}^{A}$ for some $z \in \rho(A)$ then the convergence holds for every $z \in \rho(A)$.

Proof of Theorem (6) Part (1). The discrete spectrum of $H$ can be at most $\|U\| \leq m(\mu)$ away from that of $H_{2}$. We therefore only need to consider the shift of the eigenvalues of $\mathrm{H}_{2}$ and possible production of eigenvalues from the thresholds of $H_{2}$.

Consider the eigenvalues near the origin. By standard arguments, see e.g. [25], the kernel of $H$ is spanned by linear combinations of matrix vectors composed the set $\left\{T_{j} \alpha\right\}_{j=1}^{n}$ where $\left\{T_{j}\right\}_{j=1}^{n}$ is the set of generators of symmetries of the soliton manifold, in which $\alpha$ lies. In our case there are only two symmetries: phase rotation and energy translation. The kernel is then spanned by matrix linear combinations of $\alpha$ and $\partial_{\mu} \alpha$ and thereby there exists an eigenvalue of multiplicity 2 at the origin. These eigenvalues must be result of the shift of the eigenvalues of $\mathrm{H}_{2}$ to the origin.

Now consider the possibility of eigenvalues near the threshold. Consider that by the resolvent identity, one has $R_{z}^{H}-R_{z}^{H_{2}}=R_{z}^{H_{2}} U R_{z}^{H}$. Without loss of generality, let $z$ be chosen so that $\left\|R_{z}^{H_{2}}\right\| \leq|z+\epsilon(\mu)|^{-1}$ and $|z+\epsilon(\mu)|^{-1}\|U\|<1$ for some
$\epsilon(\mu)=\mathcal{O}\left(\mu^{-\sigma}\right)$, i.e. $\epsilon(\mu)$ is due to the presence of the eigenvalues of $H_{2}$. Then

$$
\begin{equation*}
\lim _{\mu \nearrow \infty}\left\|R_{z}^{H_{2}}\right\| \leq \lim _{\mu \nearrow \infty}|z+\epsilon(\mu)|^{-1}=|z|^{-1} \tag{436}
\end{equation*}
$$

and

$$
\begin{align*}
\left\|R_{z}^{H}\right\| & =\left\|\left(1-R_{z}^{H_{2}} U\right)^{-1} R_{z}^{H_{2}}\right\| \leq\left(1-\left\|R_{z}^{H_{2}}\right\|\|U\|\right)^{-1}\left\|R_{z}^{H_{2}}\right\|  \tag{437}\\
& \leq\left[1-|z+\epsilon(\mu)|^{-1} m(\mu)\right]^{-1}|z+\epsilon(\mu)|^{-1}  \tag{438}\\
\lim _{\mu \nearrow \infty}\left\|R_{z}^{H}\right\| & \leq \lim _{\mu \nearrow \infty}\left[1-|z+\epsilon(\mu)|^{-1} m(\mu)\right]^{-1}|z+\epsilon(\mu)|^{-1}=|z|^{-1} \tag{439}
\end{align*}
$$

One may then find

$$
\begin{align*}
\lim _{\mu \nearrow \infty}\left\|R_{z}^{H}-R_{z}^{H_{2}}\right\| & =\lim _{\mu \nearrow \infty}\left\|R_{z}^{H_{2}} U R_{z}^{H}\right\| \leq \lim _{\mu \nearrow \infty}\left\|R_{z}^{H_{2}}\right\|\|U\|\left\|R_{z}^{H}\right\|  \tag{440}\\
& \leq \lim _{\mu \nearrow \infty}|z|^{-2} m(\mu)=0 \tag{441}
\end{align*}
$$

It is the case that $A$ is closed if $R_{z}^{A}$ exists and is bounded for at least one $z \in \mathbb{C}$. This is clearly the case for both $H_{2}, H$. Therefore by the principle of norm resolvent convergence it is the case that $(-\mu, 0) \cup(0, \mu) \subset \rho(H)$.

Proof of Theorem (6) Part (2). By Weyl's critereon $\sigma_{\mathrm{e}}(H)=\sigma_{\mathrm{e}}\left(H_{2}\right)$. One may explicitly construct an absolutely continuous spectral measure by expanding $R_{z}^{H}=\left(1-R_{z}^{H_{2}} U\right)^{-1} R_{z}^{H_{2}}$ as a convergent series in $U$, taking a limit $z \rightarrow \lambda \in \sigma_{\mathrm{e}}\left(H_{2}\right)$, and collecting the imaginary terms. The representation of $H v=z v$ as a coupled series of algebraic equations guarantees that each $\lambda \in \sigma_{\mathrm{e}}(H)$ has multiplicity 1 . Therefore one must have that $\sigma_{\mathrm{e}}(H)=\sigma_{\mathrm{ac}}\left(H_{2}\right)$.

## 4. Decay Estimates for $H_{2}$ and $H$

We recall without proof Lemma 3.12 from [21]:

Lemma. Let $\mathscr{B}$ be a Banach space and $\lambda_{+}>\lambda_{-}$be real constants. If $F(\lambda)$ has the properties
(1) $F \in C\left(\lambda_{-}, \lambda_{+} ; \mathscr{B}\right)$
(2) $F\left(\lambda_{-}\right)=F(\lambda)=0, \quad \lambda>\lambda_{+}$
(3) $\mathrm{d}_{\lambda} F \in L^{1}\left(\lambda_{-}+\delta, \lambda_{+} ; \mathscr{B}\right), \quad \forall \delta>0$
(4) $\mathrm{d}_{\lambda} F(\lambda)=\mathcal{O}\left(\left[\lambda-\lambda_{-}\right]^{-1} \log ^{-3}\left[\lambda-\lambda_{-}\right]\right), \quad \lambda \searrow \lambda_{-}$
(5) $\mathrm{d}_{\lambda}^{2} F(\lambda)=\mathcal{O}\left(\left[\lambda-\lambda_{-}\right]^{-2} \log ^{-2}\left[\lambda-\lambda_{-}\right]\right), \quad \lambda \searrow \lambda_{-}$
then

$$
\begin{equation*}
\int_{\lambda_{-}}^{\infty} \mathrm{d} \lambda e^{-i t \lambda} F(\lambda)=\mathcal{O}\left(t^{-1} \log ^{-2} t\right), \quad t \nearrow \infty \tag{442}
\end{equation*}
$$

in the norm of $\mathscr{B}$.

We will verify that $F(\lambda)=\delta_{\lambda}^{H_{2}}$ satisfies the desired properties for both $\lambda \geq \mu$ and $\lambda \leq \mu$.

Proof of Theorem (7). Let $\mathscr{B}=\left\{A \in \mathcal{L}(\overrightarrow{\mathscr{T}}):\|A\|_{\mathscr{B}}<\infty\right\}$ be the Banach space complete in the norm

$$
\begin{equation*}
\|A\|_{\mathscr{B}}:=\sup _{v \in \ell^{1}} \frac{\left\|W_{\kappa, \tau} A W_{\kappa, \tau} v\right\|_{1}}{\|v\|_{1}}, \tag{443}
\end{equation*}
$$

where $\overrightarrow{\mathscr{T}}$ is the natural extension of $\mathscr{T}$ to the matrix system. Let $F(\lambda)=\delta_{\lambda}^{H_{2}}$. We will verify the appropriate properties of $F(\lambda)$ for $\lambda_{-}=0$ and $\lambda_{+}=\infty$.

Let $X_{1}:=\left[\begin{array}{ll}1 & 0 \\ 0 & 0\end{array}\right], X_{2}:=\left[\begin{array}{ll}0 & 0 \\ 0 & 1\end{array}\right], g_{\lambda}:=\left[\left(1-q_{1} \mathcal{P} f_{\lambda}\right)^{2}+\left(q_{1} \pi w_{\lambda}\right)^{2}\right]^{-1}$, and

$$
\begin{align*}
& \widehat{g}_{1, \lambda}:=\left[\left(1-q_{2}^{2} \mathcal{P} f_{\lambda_{1}}^{L} f_{\lambda_{2}}^{L}\right)^{2}+\left(\pi w_{\lambda_{1}}^{L} q_{2}^{2} f_{\lambda_{2}}^{L}\right)^{2}\right]^{-1},  \tag{444}\\
& \widehat{g}_{2, \lambda}:=\left[\left(1-q_{2}^{2} f_{\lambda_{1}}^{L} \mathcal{P} f_{\lambda_{2}}^{L}\right)^{2}+\left(\pi w_{\lambda_{2}}^{L} q_{2}^{2} f_{\lambda_{1}}^{L}\right)^{2}\right]^{-1} . \tag{445}
\end{align*}
$$

Since $\psi_{z}^{L}=\left(1-q_{1} f_{z}\right)^{-1} \psi_{z}, f_{z}^{L}=\left(1-q_{1} f_{z}\right)^{-1} f_{z}$, and $\psi_{z}=f_{z} \phi_{z}+\xi_{z}$, for $\lambda \geq 0$ one has

$$
\begin{align*}
\mathcal{P} \psi_{\lambda}^{L} & =g_{\lambda}\left[\mathcal{P} f_{\lambda} \phi_{\lambda}-q_{1} \mathcal{P} f_{\lambda} \xi_{\lambda}+\xi_{\lambda}-q_{1}\left(\mathcal{P} f_{\lambda}\right)^{2} \phi_{\lambda}-q_{1}\left(\pi w_{\lambda}\right)^{2} \phi_{\lambda}\right]  \tag{446}\\
\mathcal{P} f_{\lambda}^{L} & =g_{\lambda}\left[\mathcal{P} f_{\lambda}-q_{1}\left(\mathcal{P} f_{\lambda}\right)^{2}-q_{1}\left(\pi w_{\lambda}\right)^{2}\right],  \tag{447}\\
\phi_{\lambda}^{L} & =\phi_{\lambda}+q_{1} \xi_{\lambda}, \quad w_{\lambda}^{L}=g_{\lambda} w_{\lambda} \tag{448}
\end{align*}
$$

and for $\lambda<0$ one has

$$
\begin{equation*}
\psi^{L}=\left(1-q_{1} f_{\lambda}\right)^{-1} \psi_{\lambda}, \quad f^{L}=\left(1-q_{1} f_{\lambda}\right)^{-1} f_{\lambda} \tag{449}
\end{equation*}
$$

By the method of spectral shifts one has

$$
\begin{align*}
R_{z}^{H_{2}}= & R_{z}^{H_{1}}+R_{z}^{H_{1}} q_{2} \Pi_{0} J R_{z}^{H_{2}}  \tag{450}\\
= & R_{z}^{H_{1}}+R_{z}^{H_{1}} q_{2} \Pi_{0} J\left(1-f_{z}^{H_{1}} q_{2} J\right)^{-1} R_{z}^{H_{1}} \\
= & R_{z}^{H_{1}}+R_{z}^{H_{1}} q_{2} \Pi_{0} J\left(1-q_{2}^{2} f_{z_{1}}^{L} f_{z_{2}}^{L}\right)^{-1}\left(1+f_{z}^{H_{1}} q_{2} J\right) R_{z}^{H_{1}} \\
= & R_{z}^{H_{1}}+q_{2}\left|1-q_{2}^{2} f_{z_{1}}^{L} f_{z_{2}}^{L}\right|^{-2}\left(1-q_{2}^{2} \bar{f}_{z_{1}}^{L} \bar{f}_{z_{2}}^{L}\right) R_{z}^{H_{1}}\left(J+q_{2} J f_{z}^{H_{1}} J\right) \Pi_{0} R_{z}^{H_{1}} \\
= & \left(R_{z_{1}}^{L} X_{1}-q_{2} R_{z_{2}}^{L} X_{2}\right)+\left|1-q_{2}^{2} f_{z_{1}}^{L} f_{z_{2}}^{L}\right|^{-2}\left(1-q_{2}^{2} \bar{f}_{z_{1}}^{L} \bar{f}_{z_{2}}^{L}\right) \\
& \quad \times\left(R_{z_{1}}^{L} X_{1}-R_{z_{2}}^{L} X_{2}\right)\left[J-q_{2}\left(f_{z_{1}}^{L} X_{2}-f_{z_{2}}^{L} X_{1}\right)\right] \Pi_{0}\left(R_{z_{1}}^{L} X_{1}-R_{z_{2}}^{L} X_{2}\right) \tag{455}
\end{align*}
$$

from which one finds for $\lambda \geq \mu$ :

$$
\begin{align*}
\lim _{\epsilon \searrow 0} R_{\lambda \pm i \epsilon}^{H_{2}}=[ & \left.\left(\mathcal{P}_{\lambda_{1}}^{L} X_{1}-R_{\lambda_{2}}^{L} X_{2}\right) \pm i \pi \delta_{\lambda_{1}}^{L} X_{1}\right]  \tag{456}\\
& +q_{2} \widehat{g}_{1, \lambda}\left[\left(1-q_{2}^{2} \mathcal{P} f_{\lambda_{1}}^{L} f_{\lambda_{2}}^{L}\right) \pm i \pi w_{\lambda_{1}}^{L} q_{2}^{2} f_{\lambda_{2}}^{L}\right]  \tag{457}\\
& \times\left[\left(\mathcal{P}_{z_{1}}^{L} X_{1}-R_{z_{2}}^{L} X_{2}\right) \pm i \pi \delta_{\lambda_{1}}^{L} X_{1}\right]  \tag{458}\\
& \times\left[\left(J-q_{2} \mathcal{P} f_{\lambda_{1}}^{L} X_{2}+q_{2} f_{\lambda_{2}}^{L} X_{1}\right) \mp i \pi w_{\lambda_{1}}^{L} q_{2} X_{2}\right]  \tag{459}\\
& \times \Pi_{0}\left[\left(\mathcal{P}_{\lambda_{1}}^{L} X_{1}-R_{\lambda_{2}}^{L} X_{2}\right) \pm i \pi \delta_{\lambda_{1}}^{L} X_{1}\right] \tag{460}
\end{align*}
$$

$$
\begin{align*}
=[ & \left.\left(\mathcal{P}_{\lambda_{1}}^{L} X_{1}-R_{\lambda_{2}}^{L} X_{2}\right) \pm i \pi w_{\lambda_{1}}^{L} \phi_{\lambda_{1}}^{L} \otimes \phi_{\lambda_{1}}^{L, *} X_{1}\right]  \tag{461}\\
& +q_{2} \widehat{g}_{1, \lambda}\left[\left(1-q_{2}^{2} \mathcal{P} f_{\lambda_{1}}^{L} f_{\lambda_{2}}^{L}\right) \pm i \pi w_{\lambda_{1}}^{L} q_{2}^{2} f_{\lambda_{2}}^{L}\right]  \tag{462}\\
& \times\left[\left(\mathcal{P} \psi_{\lambda_{1}}^{L} X_{1}-\psi_{\lambda_{2}}^{L} X_{2}\right) \pm i \pi w_{\lambda_{1}}^{L} \phi_{\lambda_{1}}^{L} X_{1}\right]  \tag{463}\\
& {\left[\left(J-q_{2} \mathcal{P} f_{\lambda_{1}}^{L} X_{2}+q_{2} f_{\lambda_{2}}^{L} X_{1}\right) \mp i \pi w_{\lambda_{1}}^{L} q_{2} X_{2}\right] }  \tag{464}\\
& \otimes\left[\left(\mathcal{P} \psi_{\lambda_{1}}^{L, *} X_{1}-\psi_{\lambda_{2}}^{L, *} X_{2}\right) \pm i \pi w_{\lambda_{1}}^{L} \phi_{\lambda_{1}}^{L, *} X_{1}\right] \tag{465}
\end{align*}
$$

and for $\lambda \leq \mu$ :

$$
\begin{align*}
\lim _{\epsilon \searrow 0} R_{\lambda \pm i \epsilon}^{H_{2}}=[( & \left.\left.R_{\lambda_{1}}^{L} X_{1}-\mathcal{P}_{\lambda_{2}}^{L} X_{2}\right) \mp i \pi \delta_{\lambda_{2}}^{L} X_{2}\right]  \tag{466}\\
& +q_{2} \widehat{g}_{2, \lambda}\left[\left(1-q_{2}^{2} f_{\lambda_{1}}^{L} \mathcal{P} f_{\lambda_{2}}^{L}\right) \mp i \pi w_{\lambda_{2}}^{L} q_{2}^{2} f_{\lambda_{1}}^{L}\right]  \tag{467}\\
& \times\left[\left(R_{\lambda_{1}}^{L} X_{1}-\mathcal{P}_{\lambda_{2}}^{L} X_{2}\right) \mp i \pi \delta_{\lambda_{2}}^{L} X_{2}\right]  \tag{468}\\
& {\left[\left(J-q_{2} f_{\lambda_{1}}^{L} X_{2}+q_{2} \mathcal{P} f_{\lambda_{2}}^{L} X_{1}\right) \pm i \pi w_{\lambda_{2}}^{L} q_{2} X_{1}\right] }  \tag{469}\\
& \times \Pi_{0}\left[\left(R_{\lambda_{1}}^{L} X_{1}-\mathcal{P}_{\lambda_{2}}^{L} X_{2}\right) \mp i \pi \delta_{\lambda_{2}}^{L} X_{2}\right]  \tag{470}\\
=[ & \left.\left(R_{\lambda_{1}}^{L} X_{1}-\mathcal{P}_{\lambda_{2}}^{L} X_{2}\right) \mp i \pi w_{\lambda_{2}}^{L} \phi_{\lambda_{2}}^{L} \otimes \phi_{\lambda_{2}}^{L, *} X_{2}\right]  \tag{471}\\
& +q_{2} \widehat{g}_{2, \lambda}\left[\left(1-q_{2}^{2} f_{\lambda_{1}}^{L} \mathcal{P} f_{\lambda_{2}}^{L}\right) \mp i \pi w_{\lambda_{2}}^{L} q_{2}^{2} f_{\lambda_{1}}^{L}\right]  \tag{472}\\
& \times\left[\left(\psi_{\lambda_{1}}^{L} X_{1}-\mathcal{P} \psi_{\lambda_{2}}^{L} X_{2}\right) \mp i \pi w_{\lambda_{2}}^{L} \phi_{\lambda_{2}}^{L} X_{2}\right]  \tag{473}\\
& {\left[\left(J-q_{2} f_{\lambda_{1}}^{L} X_{2}+q_{2} \mathcal{P} f_{\lambda_{2}}^{L} X_{1}\right) \pm i \pi w_{\lambda_{2}}^{L} q_{2} X_{1}\right] }  \tag{474}\\
& \otimes\left[\left(\psi_{\lambda_{1}}^{L} X_{1}-\mathcal{P} \psi_{\lambda_{2}}^{L} X_{2}\right) \mp i \pi w_{\lambda_{2}}^{L} \phi_{\lambda_{2}}^{L, *} X_{2}\right] . \tag{475}
\end{align*}
$$

One may expand the above expressions and look for the resulting imaginary piece to find $\delta_{\lambda}^{H_{2}}$ for $\lambda \geq \mu$ or $\lambda \leq \mu$. We will not do so and will analyze its properties from the unexpanded forms for simplicity instead.

We will once again use the definitions for $W_{\kappa, \tau}, \epsilon$, and the like from Chapter 1. Furthermore we will employ the spectral decay estimates of Corollary (4.1) as well as the quasi-exponential decay estimates of Theorem (4) of Chapter 02. From the latter
one can see that for $\lambda>\mu$ it is the case that $\mathrm{d}_{\lambda}^{n}\left[w_{\lambda_{1}}^{1 / 2} \psi_{\lambda_{2}}(x)\right] \in \ell^{1}$, as a function of $(x, \lambda)$, is $\ell^{1}\left(\mathbb{Z}_{+} \times[\mu, \infty), \mathbb{R}\right)$.

By considering the many definitions, there is one crucial function which strongly determines our estimates: $\mathcal{P} f_{a}=e^{a} E_{1}(a) \sim-\log (a)$ as $0<a \searrow 0$. Only powers of $\mathcal{P} f_{a}$ can be nonanalytic or unbounded. We will therefore proceed to prove the desired properties of $F(\lambda)$ by addressing the powers of $\mathcal{P} f_{a}$ alone.

We recall that by inspection one may observe that $g_{\lambda}:=\left\{\left[1-q_{1} e^{-\lambda} \mathcal{P} E_{1}(-\lambda)\right]^{2}+\right.$ $\left.\left[\pi q_{1} e^{-\lambda}\right]^{2}\right\}^{-1}$ has the properties:

$$
\begin{align*}
g_{\lambda} & =\left|g_{\lambda}\right| \leq \widehat{g}_{0}\left(q_{1}\right)<\infty, \quad \forall \lambda \in[0, \infty)  \tag{476}\\
\left|\mathrm{d}_{\lambda} g_{\lambda}\right| & \leq \widehat{g}_{0}\left(q_{1}\right) \widehat{g}_{1}\left(q_{1}, \delta\right)<\infty, \quad \forall \lambda \in[\delta, \infty)  \tag{477}\\
g_{0} & =g_{\infty}=0  \tag{478}\\
\mathrm{~d}_{\lambda} g_{\lambda} & =\mathcal{O}\left(\lambda^{-1} \log ^{-1} \lambda\right), \quad \lambda \searrow 0  \tag{479}\\
\mathrm{~d}_{\lambda}^{2} g_{\lambda} & =\mathcal{O}\left(\lambda^{-2} \log ^{-3} \lambda\right), \quad \lambda \searrow 0  \tag{480}\\
& \subset \mathcal{O}\left(\lambda^{-2} \log ^{-2} \lambda\right) \tag{481}
\end{align*}
$$

where $0<\widehat{g}_{0}\left(q_{1}\right), \widehat{g}_{1}\left(q_{1}, \delta\right)<\infty$ are constants whose other properties are not needed here. $g_{\lambda}$ is the only function of $\lambda$ involved in the definition of $F(\lambda)$ whose derivatives are unbounded in the neighborhood of the threshold $\lambda=0$ and thereby the derivatives of $g_{\lambda}$ and positive powers of $\mathcal{P} f_{\lambda}$ and its derivatives are dominant in determining the properties of the derivatives of $F(\lambda)$. We will therefore only consider the dominant factors with respect to these quantities.

Consider the contributions to the imaginary part of $\lim _{\epsilon \backslash 0} R_{\lambda \pm i \epsilon}^{H_{2}}$ for either $\lambda \geq \mu$ or $\lambda \leq \mu$. Due to the symmetry between these two ranges of $\lambda$ it is sufficient to analyze the case of $\lambda \geq \mu$ alone and we will do so exclusively in the following.
Properties (1), (2): By considering the control that the factors of $w_{\lambda_{1}}^{L}=g_{\lambda_{1}} w_{\lambda_{1}}$ impose, one may see that all possible contributions are bounded in $\lambda_{1}$ and $x$.

Property (3): The bounds for the derivatives of $\phi_{\lambda_{1}}, \xi_{\lambda_{1}}$, and $\psi_{\lambda_{2}}$ cannot present a problem with the chosen norm on $\mathscr{B}$. Exponential decay as $\lambda_{1} \nearrow \infty$ ensures that the upper bound of integration cannot be a problem. The only remaining potential issue comes from the behavior at the threshold, which is not relevant.

Property (4): The bounds for the derivatives of $\phi_{\lambda_{1}}, \xi_{\lambda_{1}}$, and $\psi_{\lambda_{2}}$ cannot present a problem with the chosen norm on $\mathscr{B}$. There will be two dominant factors. One of is $\mathrm{d}_{\lambda} g_{\lambda_{1}}=\mathcal{O}\left(\lambda_{1}^{-1} \log ^{-1} \lambda_{1}\right)$ as $\lambda_{1} \searrow 0$. The other dominant factor is of the form $g_{\lambda_{1}}^{2} \mathrm{~d}_{\lambda_{1}}\left(\mathcal{P} f_{\lambda_{1}}\right)^{2}=\mathcal{O}\left(\lambda_{1}^{-1} \log ^{-1} \lambda_{1}\right)$ as $\lambda_{1} \searrow 0$.
Property (5): The bounds for the derivatives of $\phi_{\lambda_{1}}, \xi_{\lambda_{1}}$, and $\psi_{\lambda_{2}}$ cannot present a problem with the chosen norm on $\mathscr{B}$. There will be one dominant factor, which is $g_{\lambda_{1}}\left[\mathrm{~d}_{\lambda_{1}}\left(\mathcal{P} f_{\lambda_{1}}\right)^{2}\right]^{2}=\mathcal{O}\left(\lambda^{-2} \log ^{-2} \lambda\right)$.

Proof of Theorem (8). By the Duhamel formula it is the case that

$$
\begin{align*}
-i \mathrm{~d}_{t} u= & H u=H_{2} u+U u  \tag{482}\\
\Rightarrow \quad u(t)= & e^{-i t H_{2}} u_{0}-i \int_{0}^{t} \mathrm{~d} t_{1} e^{-i\left(t-t_{1}\right) H_{2}} U u(t)  \tag{483}\\
\Rightarrow \quad W_{\kappa, \tau} u(t)= & W_{\kappa, \tau} e^{-i t H_{2}} u_{0}-i \int_{0}^{t} \mathrm{~d} t_{1} W_{\kappa, \tau} e^{-i\left(t-t_{1}\right) H_{2}} U u\left(t_{1}\right)  \tag{484}\\
\left\|W_{\kappa, \tau} u(t)\right\|_{1} \leq & \left\|W_{\kappa, \tau} e^{-i t H_{2}} u_{0}\right\|_{1}+\int_{0}^{t} \mathrm{~d} t_{1}\left\|W_{\kappa, \tau} e^{-i\left(t-t_{1}\right) H_{2}} U u\left(t_{1}\right)\right\|_{1}  \tag{485}\\
\leq & c_{0}\left\|W_{\kappa, \tau} e^{-i t H_{2}} W_{\kappa, \tau} W_{\kappa, \tau}^{-1} u_{0}\right\|_{\infty}  \tag{486}\\
& \quad+c_{1} \int_{0}^{t} \mathrm{~d} t_{1}\left\|W_{\kappa, \tau} e^{-i\left(t-t_{1}\right) H_{2}} W_{\kappa, \tau} W_{\kappa, \tau}^{-1} U W_{\kappa, \tau}^{-1} W_{\kappa, \tau} u\left(t_{1}\right)\right\|_{\infty}  \tag{487}\\
\leq & c_{0}\left\|W_{\kappa, \tau} e^{-i t H_{2}} W_{\kappa, \tau}\right\|\left\|W_{\kappa, \tau}^{-1} u_{0}\right\|_{1}  \tag{488}\\
& \quad+c_{1} \int_{0}^{t} \mathrm{~d} t_{1}\left\|W_{\kappa, \tau} e^{-i\left(t-t_{1}\right) H_{2}} W_{\kappa, \tau}\right\|\left\|W_{\kappa, \tau}^{-2} U\right\|\left\|W_{\kappa, \tau} u\left(t_{1}\right)\right\|_{1}  \tag{489}\\
\leq & c_{2}\left[\left(t+c_{3}\right) \log ^{2}\left(t+c_{3}\right)\right]^{-1}  \tag{490}\\
& \quad+c_{4} \int_{0}^{t} \mathrm{~d} t_{1}\left[\left(t-t_{1}+c_{3}\right) \log ^{2}\left(t-t_{1}+c_{3}\right)\right]^{-1}\left\|W_{\kappa, \tau} u\left(t_{1}\right)\right\|_{1} \tag{491}
\end{align*}
$$

Let $f(t):=\left\|W_{\kappa, \tau} u(t)\right\|_{1}$ and $g(t):=\left[\left(t+c_{3}\right) \log ^{2}\left(t+c_{3}\right)\right]^{-1}$. By Gronwall's Lemma one has

$$
\begin{align*}
f(t) & \leq c_{2} g(t)+c_{4} \int_{0}^{t} \mathrm{~d} t_{1} g\left(t-t_{1}\right) f\left(t_{1}\right)  \tag{492}\\
\Rightarrow \quad f(t) & \leq c_{2} g(t)+c_{2} c_{4} \int_{0}^{t} \mathrm{~d} t_{1} g\left(t_{1}\right) g\left(t-t_{1}\right) \exp \left[c_{4} \int_{t_{1}}^{t} \mathrm{~d} t_{2} g\left(t-t_{2}\right)\right] \tag{493}
\end{align*}
$$

Furthermore

$$
\begin{align*}
\exp \left[c_{4} \int_{t_{1}}^{t} \mathrm{~d} t_{2} g\left(t-t_{2}\right)\right] & =\exp \left[c_{4} \int_{0}^{t-t_{1}} \mathrm{~d} t_{2} g\left(t_{2}\right)\right]  \tag{494}\\
& \leq \exp \left[c_{4} \int_{0}^{\infty} \mathrm{d} t_{2} g\left(t_{2}\right)\right] \leq c_{5} \tag{495}
\end{align*}
$$

and

$$
\begin{align*}
\int_{0}^{t} \mathrm{~d} t_{1} g\left(t_{1}\right) g\left(t-t_{1}\right) & =2 \int_{0}^{t / 2} \mathrm{~d} t_{1} g\left(t_{1}\right) g\left(t-t_{1}\right) \leq 2 \int_{0}^{t / 2} \mathrm{~d} t_{1} g\left(t_{1}\right) g(t / 2)  \tag{496}\\
& \leq c_{6}\left[\left(t+c_{7}\right) \log ^{2}\left(t+c_{7}\right)\right]^{-1} \tag{497}
\end{align*}
$$

One therefore has

$$
\begin{align*}
f(t) & \leq c_{2} g(t)+c_{2} c_{4} c_{5} c_{6}\left[\left(t+c_{7}\right) \log ^{2}\left(t+c_{7}\right)\right]^{-1}  \tag{498}\\
& \leq c_{8}\left[\left(t+c_{9}\right) \log ^{2}\left(t+c_{9}\right)\right]^{-1}=\mathcal{O}\left(t^{-1} \log ^{-2} t\right) \tag{499}
\end{align*}
$$

as $t \nearrow \infty$. In the above $c_{j}, j=0, \ldots, 9$, are constants, the properties of which are not important.

## CHAPTER 4

## Conclusions and Conjectures

The goal for the work at this stage has been ultimately to arrive at the proof of Theorem (8). This is the most important element of the proof of asymptotic stability. One can see from the proof itself that as long as the decay estimates are integrable, and the perturbation decays faster than the square of the inverse of the weights, one can control the evolution rather easily. This will also be the case in the context of the LNLS (298) where one employs a Duhamel formula of the form

$$
\begin{equation*}
\vec{\beta}(t)=e^{-i t H} \vec{\beta}_{0}-i \int_{0}^{t} \mathrm{~d} t_{1} e^{-i\left(t-t_{1}\right) H} \vec{\gamma}(t) \tag{500}
\end{equation*}
$$

where $\vec{\beta}_{0}=\vec{\beta}(t=0)$, to specify the evolution of $\beta$. This alone, however, is not enough to determine the evolution of $\beta$ as the parameters $\mu_{1}$ and $\nu_{1}$ are time dependent. One must include separate evolution equations for these as well.

First one assumes that $(\vec{\alpha}, \vec{\beta}(0))_{\overrightarrow{\mathscr{H}}}=\left(\vec{\alpha}, \mathrm{d}_{t} \vec{\beta}(0)\right)_{\overrightarrow{\mathscr{H}}}=0$, where $(\cdot, \cdot)_{\overrightarrow{\mathscr{H}}}$ is natural the extension of the inner product of $\mathscr{H}$ to the matrix system and $\vec{\alpha}=\left[\begin{array}{c}\alpha \\ -\alpha\end{array}\right]$. This condition ensures that $\vec{\beta}$ remains in the span of the generalized eigenvectors of $H$. Then one takes the inner product of $\vec{\alpha}$ with both sides of the LNLS (298) to arrive at

$$
\begin{equation*}
0=(\alpha, \gamma) \quad \Rightarrow \quad \mathrm{d}_{t} \nu_{1}=\left(\alpha, \alpha_{\mu_{1}}\right)^{-1}\left(\alpha, \Re \gamma_{1}\right), \quad \mathrm{d}_{t} \mu_{1}=i\left(\alpha, \partial_{\mu_{1}} \alpha_{\mu_{1}}\right)^{-1}\left(\alpha, \Im \gamma_{1}\right) \tag{501}
\end{equation*}
$$

Equations (500) and (501) together constitute the modulation equations for the NLS (206), where $\mu_{1}$ and $\nu_{1}$ are the modulation parameters [31].

Conjecture 1. The soliton manifold specified by the coordinates ( $\mu, \nu$ ) with respect to $u(t)=e^{-i(-\mu t+\nu)} \alpha_{\mu}$ is asymptotically stable under perturbed evolution via the discrete NLS (206).

It is reasonable to assume that such a claim is true and it should be the case that one can prove it with an application of bootstrapping estimates. The case of the real Nonlinear Klein-Gordon equation is likely to be much harder to address.

Consider the discrete real Nonlinear Klein-Gordon equation (rNLKG) specified by

$$
\begin{equation*}
-\partial_{t}^{2} u=L_{0} u-u^{p}, \quad 1<p \in \mathbb{Z} \tag{502}
\end{equation*}
$$

It was this equation that was first studied in the context of noncommutative field theory, see e.g. [15], and the mathematical analysis of [5], [10], and [11]. There is an approach to the rNLKG which is similar to that of the methods of linearization taken with the method of modulation equations, but it is of a different character. There, the analogue of the stationary solution will be of the form $u(t)=\cos (\mu t+\nu) \alpha_{\mu}$. Due to the presence of nonlinearity this will not be a stationary solution in general. For the NLS one could interpret the stationary state as a nonlinear variant of the evolution of an eigenfunction with an associated isolated eigenvalue. For the rNLKG one is typically lead to interpret the "quasi-stationary state" as the nonlinear variant of the evolution of a resonance function with an associated embedded eigenvalue. The coupling of the "radiation" $\beta$ to the soliton will introduce an instability and lead to a resonance with a decay time. This picture was introduced and elaborated upon in the work of Soffer and Weinstein on nonlinear resonances and the nonlinear Fermi-Golden Rule [32].

Chen, Frölich, and Walcher in [5] conjecture that Equation (502) has localized metastable solutions. This leads us to the following conjecture.

Conjecture 2. Solutions of the discrete rNLKG (502) which begin close to $u(t)=\cos (\mu t+\nu) \alpha_{\mu}$ are metastable resonance functions.

We seek to address compare and contrast the proofs of these conjectures in future work.

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