# Instability, index theorem, and exponential trichotomy for Linear Hamiltonian PDEs 

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#### Abstract

Consider a general linear Hamiltonian system $\partial_{t} u=J L u$ in a Hilbert space $X$. We assume that $L: X \rightarrow X^{*}$ induces a bounded and symmetric bi-linear form $\langle L \cdot, \cdot\rangle$ on $X$, which has only finitely many negative dimensions $n^{-}(L)$. There is no restriction on the anti-self-dual operator $J: X^{*} \supset D(J) \rightarrow X$. We first obtain a structural decomposition of $X$ into the direct sum of several closed subspaces so that $L$ is blockwise diagonalized and $J L$ is of upper triangular form, where the blocks are easier to handle. Based on this structure, we first prove the linear exponential trichotomy of $e^{t J L}$. In particular, $e^{t J L}$ has at most algebraic growth in the finite co-dimensional center subspace. Next we prove an instability index theorem to relate $n^{-}(L)$ and the dimensions of generalized eigenspaces of eigenvalues of $J L$, some of which may be embedded in the continuous spectrum. This generalizes and refines previous results, where mostly $J$ was assumed to have a bounded inverse. More explicit information for the indexes with pure imaginary eigenvalues are obtained as well. Moreover, when Hamiltonian perturbations are considered, we give a sharp condition for the structural instability regarding the generation of unstable spectrum from the imaginary axis. Finally, we discuss Hamiltonian PDEs including dispersive long wave models (BBM, KDV and good Boussinesq equations), 2D Euler equation for ideal fluids, and 2D nonlinear Schrödinger equations with nonzero conditions at infinity, where our general theory applies to yield stability or instability of some coherent states.


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## CHAPTER 1

## Introduction

In this paper, we consider a general linear Hamiltonian system

$$
\begin{equation*}
\partial_{t} u=J L u, u \in X \tag{1.1}
\end{equation*}
$$

in a real Hilbert space $X$. We assume that the operator $J: X^{*} \supset D(J) \rightarrow X$ satisfies $J^{*}=-J$ and $L: X \rightarrow X^{*}$ is bounded and satisfies $L^{*}=L$. This abstract equation is motivated by the linearization of a large class of Hamiltonian PDEs at equilibria or relative equilibria. Our first goal is to understand the structural and spectral properties of (1.1), its linear stability/instability, and the persistence of these properties under small perturbations in a general setting. Secondly, the general results on (1.1) will be applied to study the linearization at some coherent states of nonlinear Hamiltonian PDEs such as the 2-dim incompressible Euler equation, generalized Bullough-Dodd equation, Gross-Pitaevskii type equation, and some long wave models like $\mathrm{KdV}, \mathrm{BBM}$, and the good Boussinesq equations.

Our main assumption is that the quadratic form $\langle L \cdot, \cdot\rangle$ admits a decomposition $X=X_{-} \oplus \operatorname{ker} L \oplus X_{+}$, such that

$$
\operatorname{dim} X_{-}=n^{-}(L)<\infty,\left.\langle L \cdot, \cdot\rangle\right|_{X_{-}}<0, \text { and }\left.\langle L \cdot, \cdot\rangle\right|_{X_{+}} \geq \delta>0
$$

An additional regularity assumption is required when $\operatorname{dim} \operatorname{ker} L=\infty$ (see (H3) in Section 2.1). We note that there is no additional restriction on the symplectic operator $J$, which can be unbounded, noninvertible, or even with infinite dimensional kernel.

* Background: stability/instability and local dynamics near an equilibrium. As our motivation for studying the linear system (1.1) is to understand the stability/instability of and the local dynamics near coherent states (steady states, traveling waves, standing waves etc.) of a nonlinear PDE, we first give a brief discussion of several standard notions of stability/instability and local dynamics. In a simple case of an ODE system

$$
x_{t}=f(x), \quad x \in \mathbf{R}^{n},
$$

the local dynamics near an equilibrium $x_{0}$, without loss of generality assuming $x_{0}=0$, is very much related to the dynamics of its linearized equation

$$
x_{t}=A x, \quad A_{n \times n}=D f(0)
$$

On the one hand, if $A$ has an unstable eigenvalue $\lambda(\operatorname{Re} \lambda>0)$, then the above linearized equation has an exponential growing solution and is therefore linearly unstable. Here, linear stability means $e^{t A}$ is uniformly bounded for all $t \geq 0$. While it is clearly linearly stable if $\operatorname{Re} \lambda<0$ for all $\lambda \in \sigma(A)$, there might be linear solutions with polynomial growth if $\operatorname{Re} \lambda \leq 0$ for all $\lambda \in \sigma(A)$, which is often referred to as the spectrally stable case. Nonlinear instability immediately follows from spectral instability for ODEs. However, it is a much more subtle issue what properties in
addition to the spectral (or even linear) stability would ensure nonlinear stability. On the other hand, assume $\sigma_{1} \subset \sigma(A)$ and $\operatorname{Re} \lambda<\alpha($ or $\operatorname{Re} \lambda>\alpha)$ for all $\lambda \in \sigma_{1}$. Let $E_{1}$ be the eigen-space of $\sigma_{1}$ which is invariant under $e^{t A}$, then we have the spectral mapping property
(SM) there exists $C>0$ s.t. $\left|e^{t A} x\right| \leq C e^{\alpha t}|x|, \forall x \in E_{1}, t \geq 0$ (or $t \leq 0$ ).
Suppose $\alpha_{+}>\alpha_{-}$and $\sigma(A)=\sigma_{+} \cup \sigma_{-}$with $\operatorname{Re} \lambda>\alpha_{+}$for all $\lambda \in \sigma_{+}$and $\operatorname{Re} \lambda<\alpha_{-}$for all $\lambda \in \sigma_{-}$. Let $E_{ \pm}$be the eigen-spaces of $\sigma_{ \pm}$, then the above spectral mapping property ( SM ) and $\alpha_{+}>\alpha_{-}$imply an exponential dichotomy of $e^{t A}$ : in the decomposition $\mathbf{R}^{n}=E_{+} \oplus E_{-}$which is invariant under $e^{t A}$, the relative minimal exponential expanding rate of $\left.e^{t A}\right|_{E_{+}}$is greater than the maximal rate of $\left.e^{t A}\right|_{E_{-}}$. For the nonlinear ODE system, the classical invariant manifold theory, based on the cornerstone of the exponential dichotomy, implies the existence of locally invariant (pseudo-)stable and unstable manifolds near 0 . They often provide more detailed dynamic structures than the mere stability/instability and also help to organize the local dynamics.

It often happens that $f(x)$ and thus $A$ depend on a small parameter $\epsilon$, so one naturally desires to understand the dynamics of the perturbed systems for $0<|\epsilon| \ll 1$ based on that of $\epsilon=0$. A system is said to be structurally stable if its dynamics does not change qualitatively under any sufficiently small perturbation. For ODEs, it is well known that the local dynamics is structurally stable if $A$ is hyperbolic, namely $\sigma(A) \cap i \mathbf{R}=\emptyset$.

The above ODE results may serve as guidelines in the study of local dynamics of PDEs near equilibria and relative equilibria while one has to keep in mind the following issues (among others):

- Sometimes it is highly non-trivial to analyze the spectra of linearized PDEs, particularly when the linear operator is not self-adjoint and has continuous spectrum. - On the eigen-space $E_{1}$ of a spectral subset $\sigma_{1}$, the above spectral mapping type property (SM) may not hold for solutions of the linearized PDEs, due to the existence of continuous spectrum of the linearized operator (see e.g. [66]).
- Regularity issues in spatial variables can cause serious complications in proving nonlinear properties (stability/instability, local invariant manifolds, etc.) based on linear ones (spectral stability/instability, exponential dichotomy, etc.). The existing systematic results are mainly for semilinear PDEs.
* Background: regarding Hamiltonian systems. On a Hilbert space $X$, a Hamiltonian system takes the form

$$
\begin{equation*}
u_{t}=J \nabla H(u) \tag{1.2}
\end{equation*}
$$

where the symplectic operator $J: X^{*} \rightarrow X$ satisfies $J^{*}=-J$ and $H: X \rightarrow \mathbf{R}$ is the Hamiltonian energy functional. In a more general setting, $J=J(u)$ may depend on $u$ or (1.2) may be posed on a symplectic manifold $M$ where $J(u): T^{*} M \rightarrow T M$. In the classical setting, the symplectic structure $\omega \in T^{*} M \otimes T^{*} M$ is a 2-form given by

$$
\omega(u)\left(U_{1}, U_{2}\right)=\left\langle J(u)^{-1} U_{1}, U_{2}\right\rangle, U_{1,2} \in T_{u} M
$$

which is required to be closed, namely $d \omega=0$. It is standard that $H$ and $\omega$ are invariant under the Hamiltonian flow associated with (1.2). Suppose $u_{*}$ is a steady state of (1.2) (possibly in an appropriate reference frame, see examples in Chapter 11), then the linearized equation at $u_{*}$ takes the form of (1.1) with $L=\nabla^{2} H\left(u_{*}\right)$. In some cases, even though the nonlinear equation is not written in a straightforward

Hamiltonian form, the linearization at an equilibrium $u_{*}$ can still be put in the Hamiltonian form (1.1), see Section 11.5 for the example of 2D Euler equation. It is standard for Hamiltonian ODEs and also proved for many Hamiltonian PDEs that the spectrum $\sigma(J L)$ is symmetric with respect to both real and imaginary axes. Therefore, either (1.1) is spectrally unstable or its spectrum must lie on the imaginary axis. Even though the latter falls into the spectral stability category, it is often subtle to obtain properties of even the linear dynamics, such as linear stability and exponential dichotomy, based on the spectral properties, particular when there is continuous spectrum. Existing results in the literature often take advantage of the conservation of $H$ or $\omega$.

The structural stability is also more subtle even for linear Hamiltonian PDEs. On the one hand, the linearized operator $J L$ associated with the linearization of Hamiltonian PDEs arising from physics and engineering usually has most of its spectrum lie on the imaginary axis. Therefore, the structural stability results based on the hyperbolicity of $J L$ are hardly applicable. On the other hand, properties of Hamiltonian systems, such as the notions of Krein signatures and the conservation of $H$ and $\omega$, provide crucial additional tools. The structural stability of linear Hamiltonian PDEs addressed in this paper is mainly related to spectral properties and linear exponential dichotomy.

For Hamiltonian PDEs, there have been some works on local nonlinear dynamics based on properties of the linearized equations. For semilinear Hamiltonian PDEs $u_{t}=J H^{\prime}(u)$ with nonlinear terms of subcritical growth, such as nonlinear Klein-Gordon equation, nonlinear Schrödinger equation, and Gross-Pitaevskii equation, local invariant manifolds can be constructed by combining ODE techniques with dispersive estimates (e. g. [4] [38] [62]). Such results for traveling wave solutions of the generalized KdV equation had also been obtained ([37]) with the help of smoothing estimates. The construction of invariant manifolds for quasilinear PDEs is more difficult, and was only done in very few cases (e. g. [57]). However, the passing from linear to nonlinear instability, which is a much weaker statement than the existence of invariant manifolds, had been done for many quasilinear PDEs (e.g. [27] [31] [36] [50] [51]). Several techniques were introduced to overcome the difficulties of loss of derivative of nonlinear terms and the growth due to the essential spectra of the linearized operators (see above references). The passing from spectral (or linear) stability to nonlinear stability is more subtle, particularly when $\langle L u, u\rangle$ is not positive definite after the symmetry reduction. When such positivity holds, the nonlinear stability can usually be proved by using the Lyapunov functional, see e.g. [29] [30] for Hamiltonian PDEs. If such positivity fails, there is currently no general approach to study the nonlinear stability based on the linear one.

Our motivation of analyzing the linearized Hamiltonian system (1.1) in such a general form is to understand the stability/instability of and the local dynamics near a coherent state $u_{*}$ of a nonlinear Hamiltonian PDE in the form of (1.2) with $L=\nabla^{2} H\left(u_{*}\right)$. We first make some comments on the hypotheses.

On $L$, the assumption $n^{-}(L)<\infty$ is equivalent to that $H(u)$ has a finite Morse index at the critical point $u_{*}$. This assumption is automatically satisfied if $u_{*}$ is constructed by minimizing $H(u)$ subject to finitely many constraints. In applications to continuum mechanics (fluids, plasmas etc.), the PDEs are often of a noncononical Hamiltonian form $u_{t}=J(u) \nabla H(u)$, with a symplectic operator
$J(u)$ depending on the solution $u$. In many cases, the linearization at an equilibrium $u_{*}$ can still be written in the Hamiltonian form (1.1) and the assumption $n^{-}(L)<$ $\infty$ is satisfied (see Section 11.5 for the example of 2D Euler equation). The uniform positivity of $L$ on $X_{+}$could be relaxed to positivity by defining a new phase space (see Chapter 10).

In the existing literature on systems in the form of (1.1), $J^{-1}: X \rightarrow X^{*}$ is mostly assumed to be a bounded operator, which is not only for technical convenience but also natural in the sense that the symplectic 2 -form $\omega$ is defined in terms of $J^{-1}$. However, it happens that $J$ does not have a bounded inverse for many important Hamiltonian PDEs such as the KdV, BBM, the good Boussinesq equations, 2D Euler equation, etc., see Chapter 11.

The goal of this paper regarding the general Hamiltonian PDE (1.1) is to study its spectral structures, linear dynamics, as well as certain structural stability properties under the assumption $n^{-}(L)<\infty$, but without any assumption on $J$ in addition to $J^{*}=-J$. Our main general results include the symmetry of the spectrum $\sigma(J L)$, an index theorem relating certain spectral properties of $J L$ to $n^{-}(L)$ which is useful for linear stability analysis, the linear exponential trichotomy of $e^{t J L}$, and the persistence of these properties for slightly perturbed Hamiltonian systems. These results are mostly achieved based on a structural decomposition of (1.1). In Chapter 11, several Hamiltonian PDEs are studied using these general results.

In the below, we briefly describe our main results and some key ideas in the proof. More details of the main theorems can be found in Chapter 2 and proofs in later chapters.

Structural decomposition. Most of the general theorems in this paper are based on careful decompositions of the phase space into closed subspaces through which $L$ and $J L$ take rather simple block forms. One of the most fundamental decomposition is given in Theorem 2.1. In this decomposition,

$$
\begin{aligned}
& J L \longleftrightarrow\left(\begin{array}{ccccccc}
0 & A_{01} & A_{02} & A_{03} & A_{04} & 0 & 0 \\
0 & A_{1} & A_{12} & A_{13} & A_{14} & 0 & 0 \\
0 & 0 & A_{2} & 0 & A_{24} & 0 & 0 \\
0 & 0 & 0 & A_{3} & A_{34} & 0 & 0 \\
0 & 0 & 0 & 0 & A_{4} & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & A_{5} & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & A_{6}
\end{array}\right), \\
& L \longleftrightarrow\left(\begin{array}{ccccccc}
0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & B_{14} & 0 & 0 \\
0 & 0 & L_{X_{2}} & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & L_{X_{3}} & 0 & 0 & 0 \\
0 & B_{14}^{*} & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & B_{56} \\
0 & 0 & 0 & 0 & 0 & B_{56}^{*} & 0
\end{array}\right),
\end{aligned}
$$

where $L$ takes an almost diagonal block form with $L_{X_{3}} \geq \delta$ for some $\delta>0$ and $J L$ takes a blockwise upper triangular form. Moreover, all the blocks of $J L$ are bounded operators except for $A_{3}$ which is anti-self-adjoint with respect to the equivalent inner product $\left\langle L_{X_{3}} \cdot, \cdot\right\rangle$ on $X_{3}$. In particular, all other diagonal blocks are matrices and therefore have only eigenvalues of finite multiplicity. The upper triangular form of
$J L$ simplifies the spectral analysis on $J L$ tremendously and plays a fundamental role in the proof of the exponential trichotomy of $e^{t J L}$, the index formula, and the structural stability/instability of (1.1).

We briefly sketch some ideas in the construction of the decomposition here under the assumption ker $L=\{0\}$, from which the decomposition in the general case follows. First, we observe that $J L$ is anti-self-adjoint in the indefinite inner product $\langle L \cdot, \cdot\rangle$. Thus, by a Pontryagin type invariant subspace Theorem for symplectic operators in an indefinite inner product space, there exists an invariant (under $J L$ ) subspace $W \subset X$, satisfying that $\left.L\right|_{W} \leq 0$ and $\operatorname{dim} W=n^{-}(L)$.

It would be highly desirable to extend $W$ to a finite dimensional invariant subspace $\tilde{W}$ such that $\left.L\right|_{\tilde{W}}$ is non-degenerate. This would yield the invariant decomposition $X=\tilde{W} \oplus \tilde{W}^{\perp L}$, where $\tilde{W}^{\perp L}$ is the orthogonal complement of $\tilde{W}$ with respect to $\langle L \cdot, \cdot\rangle$ and $\left.L\right|_{\tilde{W} \perp L}>0$. Since $\left.J L\right|_{\tilde{W} \perp L}$ is anti-self-adjoint in the equivalent inner product $\langle L \cdot, \cdot\rangle$ and $\tilde{W}$ is finite dimensional, this immediately gives the decomposition we want.

However, such an invariant decomposition $X=\tilde{W} \oplus \tilde{W}^{\perp L}$ is in general impossible since it would imply that $L$ is non-degenerate on the subspace of generalized eigenvectors of any purely imaginary eigenvalue of $J L$ (Lemma 4.1), while the counterexample in Section 8.4 shows that $L$ can be degenerate on such subspaces of embedded eigenvalues in the continuous spectra. Our proof is by a careful decomposition of the invariant spaces $W, W^{\perp L}$ and their complements.

Exponential trichotomy. Our second result is the exponential trichotomy of $e^{t J L}$ in $X$ and more regular spaces (Theorem 2.2). More precisely, we decompose $X=E^{u} \oplus E^{c} \oplus E^{s}$, such that: $E^{u, c, s}$ are invariant under $e^{t J L}$,

$$
\operatorname{dim} E^{u}=\operatorname{dim} E^{s} \leq n^{-}(L), E^{c}=\left(E^{u} \oplus E^{s}\right)^{\perp L}
$$

and $A_{5}=\left.e^{t J L}\right|_{E^{u}}\left(A_{6}=\left.e^{t J L}\right|_{E^{s}}\right)$ has exponential decay when $t<0(t>0)$ and $\left.e^{t J L}\right|_{E^{c}}$ has possible polynomial growth for all $t$ with the optimal algebraic rate explicitly given. Roughly speaking, the unstable (stable) spaces $E^{u}\left(E^{s}\right)$ are subspaces of generalized eigenvectors of the unstable (stable) eigenvalues of $J L$ and the center space $E^{c}$ corresponds to the spectra in the imaginary axis.

Such exponential trichotomy is an important step to prove nonlinear instability, and furthermore to construct local invariant (stable, unstable, center) manifolds which are crucial for a complete understanding of the local dynamics, see, for example, $[4,16,17]$. Such exponential trichotomy or dichotomy might be tricky to get due to the spectral mapping issue, that is, generally $\sigma\left(e^{t J L}\right) \subsetneq e^{t \sigma(J L)}$. So even if the spectra of $J L$ is understood, it is still a subtle issue to prove the estimates for $e^{t J L}$. In the literature, the exponential dichotomy is usually obtained either by resolvent estimates (e.g. [26]) or compact perturbations of simpler semigroups ([72] [67]). The proofs were often technical (particularly for resolvent estimates) and only worked for specific classes of problems. Our result gives the exponential trichotomy for general Hamiltonian PDEs (1.1) with $n^{-}(L)<\infty$. Moreover, the growth rates (particularly on the center space) obtained are sharp. In particular, our sharp polynomial growth rate estimate on the center space implies a stronger result than the usual spectral mapping statement. Our proof of the exponential trichotomy which is very different from traditional methods, is based on the upper triangular form of $J L$ in the decomposition given in Theorem 2.1. It can be seen that the Hamiltonian structure of (1.1) plays an important role in the proof.

Index theorems. Our third result is an index formula to relate the counting of dimensions of some eigenspaces of $J L$ to $n^{-}(L)$. Denote the sum of algebraic multiplicities of all positive eigenvalues of $J L$ by $k_{r}$ and the sum of algebraic multiplicities of eigenvalues of $J L$ in the first quadrant by $k_{c}$. Let $k_{i}^{\leq 0}$ be the total number of nonpositive dimensions $n^{\leq 0}\left(\left.L\right|_{E_{i \mu}}\right)$ of the quadratic form $\langle L \cdot, \cdot\rangle$ restricted to the subspaces $E_{i \mu}$ of generalized eigenvectors of all purely imaginary eigenvalues $i \mu \in \sigma(J L) \cap i \mathbf{R}$ of $J L$ with positive imaginary parts, and $k_{0}^{\leq 0}$ be the number of nonpositive dimensions of $\langle L \cdot, \cdot\rangle$ restricted to the generalized kernel of $J L$ modulo ker $L$. We note that, when all purely imaginary eigenvalues are semi-simple and $\langle L \cdot, \cdot\rangle$ restricted to these kernels is non-degenerate, $k_{i}^{\leq 0}$ is equal to $k_{i}^{-}$which represents the number of purely imaginary eigenvalues (with positive imaginary parts) of negative Krein signature. The situation is more complicated if the eigenvalue is not semi-simple or even embedded into the continuous spectra. In the general case, we have

$$
\begin{equation*}
k_{r}+2 k_{c}+2 k_{i}^{\leq 0}+k_{0}^{\leq 0}=n^{-}(L) . \tag{1.3}
\end{equation*}
$$

Two immediate corollaries of (1.3) are: $n^{-}(L)=k_{0}^{\leq 0}$ implies spectral stability and the oddness of $n^{-}(L)-k_{0}^{\leq 0}$ implies linear instability. Since by (1.3) all the negative directions of $\langle L \cdot, \cdot\rangle$ are associated to eigenvalues of $J L$, conceptually the continuous spectrum of $J L$ is only associated to positive directions of $\langle L \cdot, \cdot\rangle$.

There have been lots of work on similar index formulae under various settings in the literature. In the finite dimensional case where $L$ and $J L$ are matrices, such index formula readily follows from arguments in a paper of Mackay [59], although was not written explicitly there. In the past decade, there have been lots of work trying to extend it to the infinite dimensional case. In most of these papers, $J$ is assumed to have a bounded inverse ([18] [21] [44] [47]), or $\left.J\right|_{(\operatorname{ker} J)^{\perp}}$ has a bounded inverse, as in the cases of periodic waves of dispersive PDEs ([11] [13] [32] [43]). Recently, in [46] [65], the index formulae were studied for KDV type equations in the whole line for which $J=\partial_{x}$ does not have bounded inverse. Our result (1.3) gives a generalization of these results since we allow $J$ to be an arbitrary anti-self-dual operator. In particular, $\left.J\right|_{(\operatorname{ker} J)^{\perp}}$ does not need to have a bounded inverse. This is important for applications to continuum mechanics (e.g. fluids and plasmas) where $J$ usually has an infinite dimensional kernel with 0 in the essential spectrum of $J$ in some appropriate sense (see Section 11.5 for the example of 2D Euler equation).

We should also point out some differences of (1.3) with previous index formulae even in the case with bounded $J^{-1}$. In previous works on index formula, it is assumed that $\langle L \cdot, \cdot\rangle$ is non-degenerate on $(J L)^{-1}(\operatorname{ker} L) / \operatorname{ker} L$. Under this assumption, the generalized kernel of $J L$ only have Jordan blocks of length 2 and $k_{0}^{\leq 0}=n^{-}\left(\left.L\right|_{(J L)^{-1}(\operatorname{ker} L) / \operatorname{ker} L}\right)$ (see Propositions 2.7 and 2.8). In (1.3), we do not impose such non-degeneracy assumption on $\left.L\right|_{(J L)^{-1}(\operatorname{ker} L) / \text { ker } L}$ and thus the possible structures may be much richer. In the counting of (1.3), we use $k_{i}^{\leq 0}, k_{0}^{\leq 0}$, which are the total dimensions of non-positive directions of $L$ restricted on the subspaces $E_{i \mu}$ of generalized eigenvectors of purely imaginary eigenvalues $i \mu$ or zero eigenvalue (modulo ker $L$ ). Since $\langle L \cdot, \cdot\rangle$ might be degenerate on such subspace $E_{i \mu}$ of an embedded eigenvalue (see example in Section 8.4), they can not be replaced by $k_{i}^{-}, k_{0}^{-}$(i.e. the dimensions of negative directions of $L$ ) as used in the index formula of some papers (e.g. [44]). However, in Proposition 2.3, we show that if a purely
imaginary spectral point $i \mu$ is isolated, then $L$ is non-degenerate on its generalized eigenspace $E_{i \mu}$ which consists of generalized eigenvectors only. In this case, we also get an explicit formula (2.16) for $n^{-}\left(\left.L\right|_{E_{i \mu}}\right)$ by its Jordan canonical form, which is independent of the choice of the basis realizing the canonical form. This formula suggests that even for embedded eigenvalues which might be of infinite multiplicity, the number and length of nontrivial Jordan chains are bounded in terms of $n^{-}(L)$.

Moreover, even for the case where $\langle L \cdot, \cdot\rangle$ is degenerate on $E_{i \mu}$, we give a block decomposition of $J L$ and $L$ on $E_{i \mu}$ (Proposition 2.2). In this decomposition, $L$ is blockwise diagonal and $J L$ takes an upper triangular form with three diagonal blocks corresponding to the degenerate part of $L$, the simple eigenspaces and the Jordan blocks of $i \mu$ of $J L$. Furthermore, we construct a special basis for each Jordan block such that the corresponding $L$ is in an anti-diagonal form (2.15). The above decomposition of $E_{i \mu}$ yields formula (2.16) for the case where $\left.\langle L \cdot, \cdot\rangle\right|_{E_{i \mu}}$ is nondegenerate and also plays an important role on the constructive proof of Pontryagin type invariant subspace Theorem 5.1 and the proof of structural instability Theorem 2.6. To our knowledge, the formula (2.16) and the decomposition in Proposition 2.2 are new even for the finite dimensional case.

We also note that for an eigenvalue $\lambda$ with $\operatorname{Re} \lambda \neq 0,\left.\langle L \cdot, \cdot\rangle\right|_{E_{\lambda}}=0$ and by Corollary $\left.6.1\langle L \cdot, \cdot\rangle\right|_{E_{\lambda} \oplus E_{-\bar{\lambda}}}$ is non-degenerate with

$$
\begin{equation*}
n^{-}\left(\left.L\right|_{E_{\lambda} \oplus E_{-\bar{\lambda}}}\right)=\operatorname{dim} E_{\lambda} . \tag{1.4}
\end{equation*}
$$

Therefore, we get the matrix form

$$
\left.\langle L \cdot, \cdot\rangle\right|_{E_{\lambda} \oplus E_{-\bar{\lambda}}} \longleftrightarrow\left(\begin{array}{cc}
0 & A \\
A^{*} & 0
\end{array}\right)
$$

where $A$ is a nonsingular $n \times n$ matrix with $n=\operatorname{dim} E_{\lambda}$.
Now we discuss some ideas in our proof of index formula and the decomposition in Proposition 2.2 after we briefly review previous approaches for the index formulae. Like in the literature ([44] [21]), the index formula was usually proved by reducing the eigenvalue problem $J L u=\lambda u$ to a generalized eigenvalue problem $(R-z S) v=0$ (so called linear operator pencil), where $z=-\lambda^{2}$ and $R, S$ are self-adjoint operators with $\operatorname{ker} S=\{0\}$. To get such reduction it is required that $J$ has a bounded inverse and $L$ is non-degenerate on $(J L)^{-1}(\operatorname{ker} L) / \operatorname{ker} L$. Notice that the operator $S^{-1} R$ is self-adjoint in the indefinite inner product $\langle S \cdot, \cdot\rangle$. So by the Pontryagin invariant subspace theorem ([28] [18] [48] [64]) for self-adjoint operators, there is an $n^{-}(S)$-dimensional invariant (under $S^{-1} R$ ) subspace $W$ such that $\left.\langle S \cdot, \cdot\rangle\right|_{W} \leq 0$, where

$$
n^{-}(S)=n^{-}(L)-n\left(\left.L\right|_{(J L)^{-1}(\operatorname{ker} L) / \operatorname{ker} L}\right) .
$$

Going back to the original problem $J L u=\lambda u$, an index formula can be obtained by counting the negative dimensions of $L$ on the eigenspaces for real, complex and pure imaginary eigenvalues. However, it should be pointed out that the counting in some papers used the formula (1.4), for which the required non-degeneracy of $\left.L\right|_{E_{\lambda} \oplus E_{-\bar{\lambda}}}$ seemed to be assumed but not proved.

In [32] and later also in [13] [11] [43], the index formula was proved without reference to the Pontryagin invariant subspace theorem. In these papers, some conditions on $J$ and $L$ were imposed to ensure that the generalized eigenvectors of $J L$ form a complete basis of $X$. Then the index formula follows by the arguments as in the finite dimensional case ([59]). Such requirement of a complete basis is
very strong and mostly true only in some cases where the eigenvalues of $J L$ are all discrete.

Our proof of the index formula (1.3) is based on the decomposition in Theorem 2.1, where we used the Pontryagin invariant subspace theorem for the anti-self-adjoint operator $J L$ in the indefinite inner product $\langle L \cdot, \cdot\rangle$. The proof of the detailed decompositions of $J L$ and $L$ on $E_{i \mu}$ given in Proposition 2.2, particularly the construction of the special basis realizing the Jordan canonical form, is carried out in two steps. First, in the finite dimensional case, we construct a special basis of the eigenspace $E_{i \mu}$ of $J L$ to skew-diagonalize $L$ on the Jordan blocks by using an induction argument on the length of Jordan chains. Second, for the infinite dimensional case, we decompose $E_{i \mu}$ into subspaces corresponding to degenerate eigenspaces, simple non-degenerate eigenspaces and Jordan blocks. Since the Jordan block part is finite dimensional, the special basis is constructed as in the finite dimensional case.

Hamiltonian perturbations. Our fourth main result is about the persistence of exponential trichotomy and a sharp condition for the structural stability of linear Hamiltonian systems under small Hamiltonian perturbations. Consider a perturbed Hamiltonian system $u_{t}=J_{\#} L_{\#} u$ where $J_{\#}, L_{\#}$ are small perturbations of $J, L$ in the sense of (2.24). This happens when the symplectic structure or the Hamiltonian of the system depends on some parameters.

First, we show that the exponential trichotomy of $e^{t J L}$ persists under small perturbations. More precisely, we show in Theorem 2.4 that there exists a decomposition $X=E_{\#}^{u} \oplus E_{\#}^{s} \oplus E_{\#}^{c}$, satisfying that: $E_{\#}^{u, s, c}$ are invariant under $e^{t J_{\#} L_{\#}}$ and are obtained as small perturbations of $E^{u, s, c}$ in the sense that $E_{\#}^{u, s, c}=\operatorname{graph}\left(S_{\#}^{u, s, c}\right)$ where
$S_{\#}^{u}: E^{u} \rightarrow E^{s} \oplus E^{c}, \quad S_{\#}^{s}: E^{s} \rightarrow E^{u} \oplus E^{c}, \quad S_{\#}^{c}: E^{c} \rightarrow E^{s} \oplus E^{u}, \quad\left|S_{\#}^{u, s, c}\right| \leq C \epsilon$, and $\epsilon$ is roughly the size of perturbations $L_{\#}-L$ and $J_{\#}-J$ (see (2.24)). Moreover, $e^{t J_{\#} L_{\#}}$ has exponential decay on $E_{\#}^{u}$ and $E_{\#}^{s}$ in negative and positive times respectively with at most $O(\epsilon)$ loss of decay rates compared with $\left.e^{t J L}\right|_{E^{u, s}}$; on $E_{\#}^{c}, e^{t J_{\#} L_{\#}}$ has at most small exponential growth at the rate $O(\epsilon)$. We note that $\left.J_{\#} L_{\#}\right|_{E_{\#}^{c}}$ might contain eigenvalues with small real parts which are perturbed from the spectra of $J L$ in the imaginary axis and thus the small exponential growth on $e^{t J_{\#} L_{\#}}$ is the best one can get. In the perturbed decomposition $E_{\#}^{u, s, c}$, we obtain the uniform control of the growth rate and the bounds in semigroup estimates for $e^{t J_{\#} L_{\#}}$ on $E_{\#}^{u, s, c}$. Such uniform estimates of the exponential trichotomy (or dichotomy) are important for many applications of nonlinear perturbation problems, such as the modulational instability of dispersive models (see Lemma 11.2).

We briefly discuss some ideas in the proof of Theorem 2.4. The spaces $E_{\#}^{u, s}$ are constructed as the ranges of the projection operators $\tilde{P}_{\#}^{u, s}$ by the Riesz projections associated with the operator $J_{\#} L_{\#}$ in a contour enclosing $\sigma\left(\left.J L\right|_{E^{u, s}}\right)$ and $E_{\#}^{c}=\left(E_{\#}^{u, s}\right)^{\perp L_{\#}}$. The smallness assumption (2.24) is used in the resolvent estimates to show that $E_{\#}^{u, s, c}$ are indeed $O(\epsilon)$ perturbations of $E^{u, s, c}$. It is actually not so straightforward to prove the small exponential growth of $e^{t J_{\#} L_{\#}}$ on $E_{\#}^{c}$ since the perturbation term $J\left(L_{\#}-L\right)$ may be unbounded. We again use the decomposition Theorem 2.1, where in the decomposition for $J L$, only one block is infinite dimensional, with good structure, and others blocks are all bounded.

In Theorems 2.5 and 2.6, we prove that a pure imaginary eigenvalue $i \mu \neq 0$ of $J L$ is structurally stable, in the sense that the spectra of $J_{\#} L_{\#}$ near $i \mu$ stay in the imaginary axis, if and only if either $\left.L\right|_{E_{i \mu}}>0$ or $i \mu$ is isolated and $\left.L\right|_{E_{i \mu}}<0$. In particular, when $\langle L \cdot, \cdot\rangle$ is indefinite on $E_{i \mu}$ or $i \mu$ is an embedded eigenvalue and $\langle L u, u\rangle \leq 0$ for some $0 \neq u \in E_{i \mu}$, there exist perturbed operators $J L_{\#}$ with unstable eigenvalues near $i \mu$ and $\left|L_{\#}-L\right|$ being arbitrarily small. The structural stability of finite dimensional Hamiltonian systems had been well studied in the literature (see [24] [59] and references therein). It was known that (see e.g. [59]) a purely imaginary eigenvalue $i \mu \neq 0$ is structurally stable if and only if $L$ is definite on $E_{i \mu}$. As a consequence, for a family of Hamiltonian systems, the equilibrium can lose spectral stability only by the collision of purely imaginary eigenvalues of opposite Krein signatures (i.e. sign of $\langle L \cdot, \cdot\rangle$ ) . For Hamiltonian PDEs, the situation is more subtle due to the possible embedded eigenvalues in the continuous spectrum. In [28], the linearized equation at excited states of a nonlinear Schrödinger equation was studied and the structural instability was shown for an embedded simple eigenvalue with negative signature. A similar result was also obtained in [21] for semi-simple embedded eigenvalues. The assumptions in Theorems 2.5 and 2.6 are much more general and they give a sharp condition for the structural stability of nonzero pure imaginary eigenvalues of general Hamiltonian operator $J L$. In particular, in Theorem 2.6, structural instability is proved even for the case when the embedded eigenvalue is degenerate, which was not included in [28] or [21] for linearized Schrödinger equations.

In the below, we discuss some ideas in the proof of Theorems 2.5 and 2.6. In the finite dimensional case, the structural stability of an eigenvalue $i \mu$ of $J L$ with a definite energy quadratic form $\left.L\right|_{E_{i \mu}}$ can be readily seen from an argument based on Lyapunov functions. The above intuition can be used to show structural stability in Theorem 2.5 for isolated eigenvalues with definite energy quadratic forms. The proof is more subtle for embedded eigenvalues with positive energy quadratic forms. We argue via contradiction by showing that if there is a sequence of unstable eigenvalues perturbed from $i \mu$, then this leads to a non-positive direction of $\left.L\right|_{E_{i \mu}}$. In this proof, the decomposition Theorem 2.1 again plays an important role. The proof of structural instability Theorem 2.6 is divided into several cases. When $\left.L\right|_{E_{i \mu}}$ is non-degenerate and indefinite, it can be reduced to the finite dimensional case for which we can construct a perturbed matrix to have unstable eigenvalues. In particular, in the case when $E_{i \mu}$ contains a Jordan chain on which $L$ is nondegenerate, we use the special basis in Proposition 2.2 to construct a perturbed matrix with unstable eigenvalues.

The proof is more subtle for an embedded eigenvalue $i \mu$ with non-positive and possibly degenerate $\left.\langle L \cdot, \cdot\rangle\right|_{E_{i \mu}}$. First, we construct a perturbed Hamiltonian system $J \tilde{L}_{\#}$ near $J L$ such that $i \mu$ is an isolated eigenvalue of $J \tilde{L}_{\#}$ and there is a positive direction of $\left.\tilde{L}_{\#}\right|_{E_{i \mu}\left(J \tilde{L}_{\#}\right)}$. In this construction, we use the decomposition Theorem 2.1 once again along with spectral integrals. Then by Proposition $2.3,\left.\tilde{L}_{\#}\right|_{E_{i \mu}\left(J \tilde{L}_{\#}\right)}$ is non-degenerate and is indefinite by our construction. Thus it is reduced to the previously studied cases. In a rough sense, the structural instability is induced by the resonance between the embedded eigenvalue (with $\langle L \cdot, \cdot\rangle$ non-positive in the directions of some generalized eigenvectors) and the pure continuous spectra whose
spectral space has only positive directions due to the index formula (1.3).
In some applications (see e.g. Section 11.6), it is not easy to get the uniform positivity for $\left.L\right|_{X^{+}}$(i.e. assumption (H2.b)) in an obvious space $X$ and only the positivity $\left.L\right|_{X^{+}}$is available. In Theorem 2.7, we show that under some additional assumptions ((B1)-(B5) in Section 2.6), one can construct a new phase space $Y$ such that $X$ is densely embedded into $Y$; the extension $L_{Y}$ of $L$ satisfies the uniform positivity in $\|\cdot\|_{Y} ; J_{Y}: D(J) \cap Y^{*} \rightarrow Y$ is the restriction of $J$, and $\left(J_{Y}, L_{Y}, Y\right)$ satisfy the main assumptions (H1-3). Then we can apply the theorems to $\left(J_{Y}, L_{Y}, Y\right)$.

Hamiltonian PDE models. In Chapter 11 (see also Section 2.7 for a summary), we study the stabilities and related issues of various concrete Hamiltonian PDEs based on our above general theory, including: stability of solitary and periodic traveling waves of long wave models of BBM, KDV, and good Boussinesq types; the eigenvalue problem of the form $L u=\lambda u^{\prime}$ arising from the stability of solitary waves of generalized Bullough-Dodd equation; modulational instability of periodic traveling waves; stability of steady flows of 2D Euler equations; traveling waves of 2D nonlinear Schrödinger equations with nonzero condition at infinity.

This paper is organized as follows. In Chapter 2, we give the precise set-up and list the main general results more precisely with some comments, where the readers are directed to the corresponding subsequent sections for detailed proofs. For some readers, who would like to see the general results but do not desire to get into the technical details of the proofs, it is possibly sufficient to read Sections 2.1-2.6 only. The stability analysis of various Hamiltonian PDEs are outlined in Section 2.7. The proofs of the main general results are given in Chapters 3 to 10 . Chapter 3 studies some basic properties of linear Hamiltonian systems. Chapter 4 is about the finite dimensional Hamiltonian systems. In particular, the special basis in Proposition 2.2 is constructed. Chapter 5 is about the Pontryagin type invariant subspace Theorem for anti-self-adjoint operators in an indefinite inner product space. Two proofs are given. One is by the fixed point argument as found in the literature ([18] [25] [48]), which provides the existence of an invariant Pontryagin subspace abstractly. The second one in separable Hilbert spaces is via Galerkin approximation which also yields an explicit construction of a maximally non-positive invariant subspace. Chapter 6 is to prove decomposition Theorem 2.1 which plays a crucial role in the proof of most of the main results. Chapter 7 contains the proof of the exponential trichotomy of $e^{t J L}$. In Chapter 8 , the index theorem is proved. Besides, the structures of the generalized eigenspaces are studied and more explicit formula for the indexes $k_{i}^{\leq 0}, k_{0}^{\leq 0}$, etc. are proved. The non-degeneracy of $\left.L\right|_{E_{i \mu}}$ for any isolated spectral point $i \mu$ is also proved there. In Chapter 9, we prove the persistence of the exponential trichotomy and the structural stability/instability Theorems. In Chapter 10, we prove that the uniform positivity assumption (H2.b) can be relaxed under some assumptions. We study the stability and related issues of various Hamiltonian PDEs in Chapter 11. In the Appendix, we prove some functional analysis facts used throughout the paper, including some basic decompositions of the phase space, the well-posedness of the linear Hamiltonian system, and the standard complexification procedure.

## CHAPTER 2

## Main results

In this chapter, we give details of the main results described in the introduction. The detailed proofs are left for later sections.

## A remark on notations: Throughout the paper, given a densely defined linear

 operator $T$ from a Banach space $X$ to a Banach space $Y$ we will always use $T^{*}$ to denote its dual operator from a subspace of $Y^{*}$ to $X^{*}$. It would never mean the adjoint operator even if $X=Y$ is a Hilbert space. Given a Hilbert space $X$ and a linear operator $L: X \rightarrow X^{*}$, since $L^{*}:\left(X^{*}\right)^{*}=X \rightarrow X^{*}$, it is legitimate to compare whether $L=L^{*}$.
### 2.1. Set-up

Consider a linear Hamiltonian system

$$
\begin{equation*}
\partial_{t} u=J L u, \quad u \in X \tag{2.1}
\end{equation*}
$$

where $X$ is a real Hilbert space. Let $(\cdot, \cdot)$ denote the inner product on $X$ and $\langle\cdot, \cdot\rangle$ the dual bracket between $X^{*}$ and $X$. We make the following assumptions:
(H1) $J: X^{*} \supset D(J) \rightarrow X$ is anti-self-dual, in the sense $J^{*}=-J$.
(H2) The operator $L: X \rightarrow X^{*}$ is bounded and symmetric (i.e. $L^{*}=L$ ) such that $\langle L u, v\rangle$ is a bounded symmetric bilinear form on $X$. Moreover, there exists a decomposition of $X$ into the direct sum of three closed subspaces

$$
X=X_{-} \oplus \operatorname{ker} L \oplus X_{+}, \quad n^{-}(L) \triangleq \operatorname{dim} X_{-}<\infty
$$

satisfying
(H2.a) $\langle L u, u\rangle<0$ for all $u \in X_{-} \backslash\{0\}$;
(H2.b) there exists $\delta>0$ such that

$$
\langle L u, u\rangle \geq \delta\|u\|^{2}, \text { for any } u \in X_{+}
$$

(H3) The above $X_{ \pm}$satisfy

$$
\operatorname{ker} i_{X_{+} \oplus X_{-}}^{*}=\left\{f \in X^{*} \mid\langle f, u\rangle=0, \forall u \in X_{-} \oplus X_{+}\right\} \subset D(J)
$$

where $i_{X_{+} \oplus X_{-}}^{*}: X^{*} \rightarrow\left(X_{+} \oplus X_{-}\right)^{*}$ is the dual operator of the embedding $i_{X_{+} \oplus X_{-}}$.

REMARK 2.1. If in addition we assume

$$
\begin{equation*}
\operatorname{ker} i_{(\operatorname{ker} L)^{\perp}}^{*}=\left\{f \in X^{*} \mid\langle f, u\rangle=0, \forall u \in(\operatorname{ker} L)^{\perp}\right\} \subset D(J) \tag{2.2}
\end{equation*}
$$

where

$$
\begin{equation*}
(\operatorname{ker} L)^{\perp}=\{u \in X \mid(u, v)=0, \forall v \in \operatorname{ker} L\} \tag{2.3}
\end{equation*}
$$

it is possible to choose $X_{ \pm} \subset(\operatorname{ker} L)^{\perp}$. See Lemma 12.4 and Remark 12.4.

Regarding the operator $L$, what often matters more is its associated symmetric quadratic form $\langle L u, v\rangle, u, v \in X$, (or the Hermitian symmetric form after the complexification). We say a bounded symmetric quadratic form $B(u, v)$ is nondegenerate if

$$
\begin{equation*}
\inf _{v \neq 0} \sup _{u \neq 0} \frac{|B(u, v)|}{\|u\|\|v\|}>0 \tag{2.4}
\end{equation*}
$$

or equivalently, $v \rightarrow f=B(\cdot, v) \in X^{*}$ defines an isomorphism from $X$ to $X^{*}$ (or a complex conjugate (sometimes called anti-linear) isomorphism - satisfying $a v \rightarrow \bar{a} f$ for any $a \in \mathbf{C}-$ after the complexification). Under assumptions (H1-3), $\langle L u, v\rangle$ is non-degenerate if and only if ker $L=\{0\}$ (see Lemma 12.2).

REMARK 2.2. It is worth pointing out that $n^{-}(L)=\operatorname{dim} X_{-}$is actually the maximal dimension of subspaces where $\langle L \cdot, \cdot\rangle<0$, see Lemma 12.1. Thus $n^{-}(L)$ is the Morse index of $L$.

By Riesz Representation Theorem, there exists a unique bounded symmetric linear operator $\mathbb{L}: X \rightarrow X$ such that $(\mathbb{L} u, v)=\langle L u, v\rangle$. Let $\Pi_{\lambda}, \lambda \in \mathbf{R}$, denote the orthogonal spectral projection operator from $X$ to the closed subspace corresponding to the spectral subset $\sigma(\mathbb{L}) \cap(-\infty, \lambda]$. From the standard spectral theory of selfadjoint operators, assumption (H2) is equivalent to that there exists $\delta^{\prime}>0$ such that
i.) $\sigma(\mathbb{L}) \cap\left[-\delta^{\prime}, \delta^{\prime}\right] \subset\{0\}$, which is equivalent to the closeness of $R(L)$, and
ii.) $\operatorname{dim}\left(\Pi_{-\delta^{\prime}} X\right)<\infty$.

The subspaces

$$
X_{-}=\Pi_{-\frac{\delta^{\prime}}{2}} X \quad X_{+}=\left(I-\Pi_{\frac{\delta^{\prime}}{2}}\right) X,
$$

along with $\operatorname{ker} L$ lead to a decomposition of $X$ orthogonal with respect to both $(\cdot, \cdot)$ and $\langle L \cdot, \cdot\rangle$, satisfying (H2).

REmARK 2.3. We would like to point out that (H3) is automatically satisfied if $\operatorname{dim} \operatorname{ker} L<\infty$. In fact in this case,
$\operatorname{dim} \operatorname{ker} i_{X_{+} \oplus X_{-}}^{*}=\operatorname{dim}\left\{f \in X^{*} \mid\langle f, u\rangle=0, \forall u \in X_{-} \oplus X_{+}\right\}=\operatorname{dim} \operatorname{ker} L<\infty$.
Let $\left\{f_{1}, \ldots, f_{k}\right\}$ be a basis of $\operatorname{ker} i_{X_{+} \oplus X_{-}}^{*} . A s D(J)$ is dense in $X^{*}$, one may take $g_{j} \in D(J)$ sufficiently close to $f_{j}, j=1, \ldots, k$. Let

$$
X_{1}=\left\{u \in X \mid\left\langle g_{j}, u\right\rangle=0, \forall j=1, \ldots, k\right\}
$$

Since $X_{1}$ is close to $X_{+} \oplus X_{-}$, it is easy to show that there exist closed subspaces $X_{1 \pm} \subset X_{1}$ satisfying (H2) and $X_{1}=X_{1+} \oplus X_{1-}$ 。

In fact, if we had treated $L$ and $J$ as operators from $X$ to $X$ through the Riesz Representation Theorem and $X_{ \pm}$happen to be given as in Remark 2.2 then (H3) would take the form $\operatorname{ker} L \subset D(J)$.

Assumption (H3) does ensure that $J L$ is densely defined, see Lemma 12.5.
Remark 2.4. Assumption (H2.b) requires that the quadratic form $\langle L u, u\rangle$ has a uniform positive lower bound on $X_{+}$. This corresponds to that 0 is an isolated eigenvalue of $\mathbb{L}$ defined in Remark 2.2, which also implies that $R(L)$ is closed and $R(L)=\left\{\gamma \in X^{*} \mid\langle\gamma, u\rangle=0, \forall u \in \operatorname{ker} L\right\}$.

For some PDE systems, (H2.b) may not hold or be hard to verify, see, e.g. Section 11.6. In Section 2.6, we consider a framework where assumption (H2.b) for the uniform positivity of $\left.L\right|_{X_{+}}$is weakened to the positivity of $\left.L\right|_{X_{+}}$, if some
additional and more detailed structures are present. In that situation, we construct a new phase space $Y \supset X$ and extend the operators $L$ and $J$ to $Y$ accordingly so that (H1-3) are satisfied.

### 2.2. Structural decomposition

Our first main result is to construct a decomposition of the phase space $X$ which helps understanding both structures of $J L$ and $L$ simultaneously.

Theorem 2.1. Assume (H1-H3). There exist closed subspaces $X_{j}, j=1, \ldots, 6$, and $X_{0}=\operatorname{ker} L$ such that
(1) $X=\oplus_{j=0}^{6} X_{j}, X_{j} \subset \cap_{k=1}^{\infty} D\left((J L)^{k}\right), j \neq 3$, and
$\operatorname{dim} X_{1}=\operatorname{dim} X_{4}, \operatorname{dim} X_{5}=\operatorname{dim} X_{6}, \operatorname{dim} X_{1}+\operatorname{dim} X_{2}+\operatorname{dim} X_{5}=n^{-}(L) ;$
(2) $J L$ and $L$ take the following forms in this decomposition

$$
\begin{gather*}
J L \longleftrightarrow\left(\begin{array}{ccccccc}
0 & A_{01} & A_{02} & A_{03} & A_{04} & 0 & 0 \\
0 & A_{1} & A_{12} & A_{13} & A_{14} & 0 & 0 \\
0 & 0 & A_{2} & 0 & A_{24} & 0 & 0 \\
0 & 0 & 0 & A_{3} & A_{34} & 0 & 0 \\
0 & 0 & 0 & 0 & A_{4} & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & A_{5} & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & A_{6}
\end{array}\right),  \tag{2.5}\\
L \longleftrightarrow\left(\begin{array}{ccccccc}
0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & B_{14} & 0 & 0 \\
0 & 0 & L_{X_{2}} & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & L_{X_{3}} & 0 & 0 & 0 \\
0 & B_{14}^{*} & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & B_{56} \\
0 & 0 & 0 & 0 & 0 & B_{56}^{*} & 0
\end{array}\right) .
\end{gather*}
$$

(3) $B_{14}: X_{4} \rightarrow X_{1}^{*}$ and $B_{56}: X_{6} \rightarrow X_{5}^{*}$ are isomorphisms and there exists $\delta>0$ satisfying $\mp\left\langle L_{X_{2,3}} u, u\right\rangle \geq \delta\|u\|^{2}$, for all $u \in X_{2,3}$;
(4) all blocks of $J L$ are bounded operators except $A_{3}$, where $A_{03}$ and $A_{13}$ are understood as their natural extensions defined on $X_{3}$;
(5) $A_{2,3}$ are anti-self-adjoint with respect to the equivalent inner product $\mp\left\langle L_{X_{2,3}} \cdot, \cdot\right\rangle$ on $X_{2,3}$;
(6) the spectra $\sigma\left(A_{j}\right) \subset i \mathbf{R}, j=1,2,3,4, \pm \operatorname{Re} \lambda>0$ for all $\lambda \in \sigma\left(A_{5,6}\right)$, and $\sigma\left(A_{5}\right)=-\sigma\left(A_{6}\right)$;
(7) $n^{-}\left(\left.L\right|_{X_{5} \oplus X_{6}}\right)=\operatorname{dim} X_{5}$ and $n^{-}\left(\left.L\right|_{X_{1} \oplus X_{4}}\right)=\operatorname{dim} X_{1}$.
(8) $(u, v)=0$ for all $u \in X_{1} \oplus X_{2} \oplus X_{3} \oplus X_{4}$ and $v \in \operatorname{ker} L$.

Through straightforward calculations, one may naturally rewrite the operator $J$ and obtain additional relations among those blocks $A_{j k}$ using $J^{*}=-J$.

Corollary 2.1. Let $P_{j}, j=0, \ldots, 6$ be the projections associated to the decomposition in Theorem 2.1 and $\tilde{X}_{j}^{*}=P_{j}^{*} X_{j}^{*} \subset X^{*}$. In the decomposition
$X^{*}=\Sigma_{j=0}^{6} \tilde{X}_{j}^{*}, J$ has the block form

$$
J \longleftrightarrow\left(\begin{array}{ccccccc}
J_{00} & J_{01} & J_{02} & J_{03} & J_{04} & 0 & 0 \\
J_{10} & J_{11} & J_{12} & J_{13} & J_{14} & 0 & 0 \\
J_{20} & J_{21} & J_{22} & 0 & 0 & 0 & 0 \\
J_{30} & J_{31} & 0 & J_{33} & 0 & 0 & 0 \\
J_{40} & J_{41} & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & J_{56} \\
0 & 0 & 0 & 0 & 0 & J_{65} & 0
\end{array}\right) .
$$

where the blocks, except $J_{00}$, are given by

$$
\begin{aligned}
& -J_{10}^{*}=J_{01}=A_{04} B_{14}^{-1}, \quad-J_{20}^{*}=J_{02}=A_{02} L_{2}^{-1} \\
& -J_{30}^{*}=J_{03}=A_{03} L_{X_{3}}^{-1}, \quad-J_{40}^{*}=J_{04}=A_{01}\left(B_{14}^{*}\right)^{-1} \\
& J_{11}=A_{14} B_{14}^{-1}, \quad J_{12}=A_{12} L_{X_{2}}^{-1}, \quad J_{13}=A_{13} L_{X_{3}}^{-1}, \quad J_{14}=A_{1}\left(B_{14}^{*}\right)^{-1} \\
& J_{21}=A_{24} B_{14}^{-1}, \quad J_{22}=A_{2} L_{X}^{-1}, \quad J_{31}=A_{34} B_{14}^{-1}, \quad J_{33}=A_{3} L_{X_{3}}^{-1} \\
& J_{41}=A_{4} B_{14}^{-1}, \quad J_{56}=A_{5}\left(B_{56}^{*}\right)^{-1}, \quad J_{65}=A_{6} B_{56}^{-1}
\end{aligned}
$$

Due to $J^{*}+J=0$, we also have $L_{X_{j}} A_{j}+A_{j}^{*} L_{X_{j}}=0, j=2,3$, and

$$
B_{14}^{*} A_{14}+A_{14}^{*} B_{14}=0, \quad L_{X_{2}} A_{24}+A_{12}^{*} B_{14}=0, \quad L_{X_{3}} A_{34}+A_{13}^{*} B_{14}=0
$$

$$
B_{14} A_{4}+A_{1}^{*} B_{14}=0, \quad B_{56} A_{6}+A_{5}^{*} B_{56}=0
$$

REmARK 2.5. From the corollary, we have the following observations.
(i) $A_{4}$ and $-A_{1}^{*}$ are similar through $B_{14}$ and thus have the same spectrum, contained in $i \mathbf{R}$ and symmetric about the real axis. This in turn implies that $\sigma\left(A_{1}\right)=\sigma\left(A_{4}\right)$.
(ii) $A_{24}$ and $A_{34}$ can be determined by other blocks

$$
A_{24}=-L_{X_{2}}^{-1} A_{12}^{*} B_{14}, \quad A_{34}=-L_{X_{3}}^{-1} A_{13}^{*} B_{14} .
$$

Consequently,

$$
J_{21}=-L_{X_{2}}^{-1} A_{12}^{*}, \quad J_{31}=-L_{X_{3}}^{-1} A_{13}^{*}
$$

The proof of Theorem 2.1 is given in Chapter 6, largely based on the Pontryagin invariant subspace theorem 5.1. Theorem 2.1 decomposes the closed operator $J L$ into an upper triangular block form, all of which are bounded except for one block anti-self-adjoint with respective to an equivalent norm. This decomposition plays a fundamental role in proving the linear evolution estimates, the index theorem, the spectral analysis, and the perturbation analysis.

### 2.3. Exponential Trichotomy

One of our main results is the exponential trichotomy of the semigroup $e^{t J L}$ on $X$ and more regular spaces, to be proved in Chapter 7. Such linear estimates are important for studying nonlinear dynamics, particularly, the construction of invariant manifolds for nonlinear Hamiltonian PDEs.

THEOREM 2.2. Under assumptions (H1)-(H3), JL generates a $C^{0}$ group e $e^{t J L}$ of bounded linear operators on $X$ and there exists a decomposition

$$
X=E^{u} \oplus E^{c} \oplus E^{s}, \quad \operatorname{dim} E^{u}=\operatorname{dim} E^{s} \leq n^{-}(L)
$$

satisfying:
i) $E^{c}$ and $E^{u}, E^{s} \subset D(J L)$ are invariant under $e^{t J L} ;$ Here, $E^{u}=X_{5}, E^{s}=X_{6}$
are the unstable and stable spaces defined in Theorem 2.1, and the center space $E^{c}$ is defined by

$$
E^{c}=\left\{u \in X \mid\langle L u, v\rangle=0, \forall v \in E^{s} \oplus E^{u}\right\}=\oplus_{j=0}^{4} X_{j}
$$

ii) $\langle L \cdot, \cdot\rangle$ completely vanishes on $E^{u, s}$, but is non-degenerate on $E^{u} \oplus E^{s}$;
iii) let $\lambda_{u}=\min \{\operatorname{Re} \lambda \mid \lambda \in \sigma(J L)$, Re $\lambda>0\}$, there exist $M>0$ and an integer $k_{0} \geq 0$, such that

$$
\begin{gather*}
\left|e^{t J L}\right|_{E^{s}} \mid \leq M\left(1+t^{\operatorname{dim} E^{s}-1}\right) e^{-\lambda_{u} t}, \quad \forall t \geq 0  \tag{2.7}\\
\left|e^{t J L}\right|_{E^{u}} \mid \leq M\left(1+|t|^{\operatorname{dim} E^{u}-1}\right) e^{\lambda_{u} t}, \quad \forall t \leq 0 \\
\left|e^{t J L}\right|_{E^{c}} \mid \leq M\left(1+|t|^{k_{0}}\right), \quad \forall t \in \mathbf{R} \tag{2.8}
\end{gather*}
$$

and

$$
k_{0} \leq 1+2\left(n^{-}(L)-\operatorname{dim} E^{u}\right)
$$

Moreover, for $k \geq 1$, define the space $X^{k} \subset X$ to be

$$
X^{k}=D\left((J L)^{k}\right)=\left\{u \in X \mid(J L)^{n} u \in X, n=1, \cdots, k .\right\}
$$

and

$$
\begin{equation*}
\|u\|_{X^{k}}=\|u\|+\|J L u\|+\cdots+\left\|(J L)^{k} u\right\| \tag{2.9}
\end{equation*}
$$

Assume $E^{u, s} \subset X^{k}$, then the exponential trichotomy for $X^{k}$ holds true: $X^{k}$ is decomposed as a direct sum

$$
X^{k}=E^{u} \oplus E_{k}^{c} \oplus E^{s}, E_{k}^{c}=E^{c} \cap X^{k}
$$

and the estimates (2.7) and (2.8) still hold in the norm $X^{k}$.
An immediate corollary of the theorem is that there are only finitely many eigenvalues of $J L$ outside the imaginary axis in the complex plane.

REMARK 2.6. The above growth estimates is optimal as one may easily construct finite dimensional examples which achieve upper bounds in the estimates.

REMARK 2.7. Naturally, the above invariant decomposition and exponential trichotomy are based on the spectral decomposition of JL. The unstable/stable subspaces $E^{u, s}$ are the eigenspaces of the stable/unstable spectrum, which have finite total dimensions. Therefore, it is easy to obtain the exponential decay estimates of $\left.e^{t J L}\right|_{E^{u, s}}$. While $E^{c}$ is the eigenspace of the spectrum residing on the imaginary axis, the growth estimate of $\left.e^{t J L}\right|_{E^{c}}$ is far from obvious as the spectral mapping is often a complicated issue especially when continuous spectra is involved. Normally some sub-exponential growth estimates, like in the form of

$$
\forall \epsilon>0, \exists C>0 \Longrightarrow\left|e^{t J L}\right|_{E^{c}} \mid \leq C e^{\epsilon|t|}, \forall t \in \mathbf{R}
$$

are already sufficient for some nonlinear local analysis. Our above polynomial growth estimate on $\left.e^{t J L}\right|_{E^{c}}$ with uniform bound on the degree of the polynomial based on $\operatorname{dim} X_{-}$is a much stronger statement.

REMARK 2.8. Often the invariant subspaces $E^{u, s, c}$ are defined via spectral decompositions where the L-orthogonality between $E^{s} \oplus E^{u}$ and $E^{c}$ is not immediately clear. In fact, this is a special case of more general L-orthogonality property. See Lemma 6.2 and Corollary 6.2.

### 2.4. Index Theorems and spectral properties

Roughly our next main result is on the relationship between the number of negative directions of $L$ (the Morse index) and the dimensions of various eigenspaces of $J L$, which may have some implications on $\operatorname{dim} E^{u, s}$ and thus the stability/instability of the group $e^{t J L}$.

We first introduce some notations. Given any subspace $S \subset X$, denote $n^{-}\left(\left.L\right|_{S}\right)$ and $n^{\leq 0}\left(\left.L\right|_{S}\right)$ as the maximal negative and non-positive dimensions of $\langle L u, u\rangle$ restricted to $S$, respectively. Clearly, $n^{-}\left(\left.L\right|_{s}\right) \leq n^{-}(L)<\infty$.

In order to state and prove our results on the index theorems, we will work with the standard complexified spaces, operators, and quadratic forms, see Appendix (Chapter 12) for details.

For any eigenvalue $\lambda$ of $J L$ let $E_{\lambda}$ be the generalized eigenspace, that is,

$$
E_{\lambda}=\left\{u \in X \mid(J L-\lambda I)^{k} u=0, \text { for some integer } k \geq 1\right\}
$$

REMARK 2.9. As $J L$ generates a $C^{0}$ semigroup (Proposition 12.1), $(J L-\lambda)^{k}$ is a densely defined closed operator (see [33]) and thus $E_{\lambda}$ is indeed a closed subspace. It will turn out that $E_{\lambda}=\operatorname{ker}(J L-\lambda I)^{2 n^{-}(L)+1}$ for any eigenvalue $\lambda$. See Theorem 2.3 for $\lambda \notin i \mathbf{R}$ and Proposition 2.1 for more details.

Let $k_{r}$ be the sum of algebraic multiplicities of positive eigenvalues of $J L$ and $k_{c}$ be the sum of algebraic multiplicities of eigenvalues of $J L$ in the first quadrant (i.e. both real and imaginary parts are positive). Namely,

$$
\begin{equation*}
k_{r}=\sum_{\lambda>0} \operatorname{dim} E_{\lambda}, \quad k_{c}=\sum_{\operatorname{Re} \lambda, \operatorname{Im} \lambda>0} \operatorname{dim} E_{\lambda} . \tag{2.10}
\end{equation*}
$$

For any purely imaginary eigenvalue $i \mu\left(0 \neq \mu \in \mathbf{R}^{+}\right)$of $J L$, let

$$
\begin{equation*}
k^{\leq 0}(i \mu)=n^{\leq 0}\left(\left.L\right|_{E_{i \mu}}\right), \quad k_{i}^{\leq 0}=\sum_{0 \neq \mu \in \mathbf{R}^{+}} k^{\leq 0}(i \mu) . \tag{2.11}
\end{equation*}
$$

The index counting on $E_{0}$ is slightly more subtle due to the possible presence of nontrivial ker $L \subset E_{0}$. Observe that, for any subspace $S \subset X, L$ induces a quadratic form $\langle L \cdot, \cdot\rangle$ on the quotient space $S /(\operatorname{ker} L \cap S)$. As ker $L \subset E_{0}$, define

$$
\begin{equation*}
k_{0}^{\leq 0}=n^{\leq 0}\left(\left.\langle L \cdot, \cdot\rangle\right|_{E_{0} / \operatorname{ker} L}\right) . \tag{2.12}
\end{equation*}
$$

Equivalently, let $\tilde{E}_{0} \subset E_{0}$ be any subspace satisfying $E_{0}=\operatorname{ker} L \oplus \tilde{E}_{0}$. Define

$$
k_{0}^{\leq 0}=n^{\leq 0}\left(\left.L\right|_{\tilde{E}_{0}}\right)
$$

It is easy to see that $k_{0}^{\leq 0}$ is independent of the choice of $\tilde{E}_{0}$. We have the following index formula which is proved in Section 8.1.

Theorem 2.3. Assume (H1)-(H3), we have
(i) If $\lambda \in \sigma(J L)$, then $\pm \lambda, \pm \bar{\lambda} \in \sigma(J L)$.
(ii) If $\lambda$ is an eigenvalue of $J L$, then $\pm \lambda, \pm \bar{\lambda}$ are all eigenvalues of JL. Moreover, for any integer $k>0$,

$$
\operatorname{dim} \operatorname{ker}(J L \pm \lambda)^{k}=\operatorname{dim} \operatorname{ker}(J L \pm \bar{\lambda})^{k}
$$

(iii) The indices satisfy

$$
\begin{equation*}
k_{r}+2 k_{c}+2 k_{i}^{\leq 0}+k_{0}^{\leq 0}=n^{-}(L) . \tag{2.13}
\end{equation*}
$$

Combining Theorems 2.2 and 2.3, we have the following corollary.

Corollary 2.2. (i) If $k_{0}^{\leq 0}=n^{-}(L)$, then (2.1) is spectrally stable. That is, there exists no exponentially unstable solution of (2.1).
(ii) If $n^{-}(L)-k_{0}^{\leq 0}$ is odd, then there exists a positive eigenvalue of (2.1), that is, $k_{r}>0$. In particular, if $n^{-}(L)-k_{0}^{\leq 0}=1$, then $k_{r}=1$ and $k_{c}=k_{i}^{\leq 0}=0$, that is, (2.1) has exactly one pair of stable and unstable simple eigenvalues.

REMARK 2.10. The formula (2.13) might seem more intuitive if those above $k^{\leq 0}$ had been replaced by $k^{-}$. In fact such an index formula with $k^{-}$instead of $k^{\leq 0}$ is true only if the quadratic form $\langle L u, v\rangle$ is non-degenerate on all $E_{i \mu}, \mu \in \mathbf{R}^{+}$and $\tilde{E}_{0}$, which would imply $n^{-}\left(\left.L\right|_{E_{i \mu}}\right)=n^{\leq 0}\left(\left.L\right|_{E_{i \mu}}\right)$. However, the degeneracy is indeed possible and the correct choice has to be $k \leq 0$. Such an example is given in Section 8.4.

Even though we can not claim $\operatorname{dim} E_{i \mu}<\infty$ for an eigenvalue $i \mu \in i \mathbf{R}$ which might be embedded in the continuous spectrum, in fact $E_{i \mu}$ is spanned by eigenvectors along with finitely many generalized eigenvectors, except for $\mu=0$. More precisely, we prove the following two propositions in Lemma 3.5 and Section 8.2.

Proposition 2.1. Assume (H1)-(H3). For any $\mu \in \sigma(J L) \cap \mathbf{R} \backslash\{0\}$, it holds

$$
E_{i \mu}=\operatorname{ker}(J L-i \mu)^{2 k^{\leq 0}(i \mu)+1}, \quad \operatorname{dim}\left((J L-i \mu) E_{i \mu}\right) \leq 2 k^{\leq 0}(i \mu)
$$

Moreover,

$$
E_{0}=\operatorname{ker}(J L)^{2 k_{0}^{\leq 0}+2}, \quad \operatorname{dim}\left((J L)^{2} E_{0}\right) \leq 2 k_{0}^{\leq 0}
$$

The above proposition does not hold if $(J L)^{2} E_{0}$ is replaced by $J L E_{0}$ as in the case of $\mu \neq 0$. See an example in Remark 8.2 in Section 8.2.

For $\mu \in \mathbf{R}$, Theorem 2.3 and Proposition 2.1 mean that, in addition to eigenvectors, $\left.J L\right|_{E_{i \mu}}$ has only finitely many nontrivial Jordan blocks with the total dimensions bounded in term of $n^{-}(L)$. The number and the lengths of nontrivial Jordan chains of $\left.J L\right|_{E_{i \mu}}$ are independent of the choice of the basis realizing the Jordan canonical form. Intuitively if a basis consisting of generalized eigenvectors simultaneously diagonalizes the quadratic form $\langle L u, u\rangle$ and realizes the Jordan canonical form of $J L$, it would greatly help us to understand the structure of (2.1). However, usually this is not possible. Instead, we find a 'good' basis for the Jordan canonical form of $J L$ which also 'almost' diagonalizes the quadratic form $L$. To our best knowledge, we are not aware of such a result even in finite dimensions.

Proposition 2.2. Assume (H1)-(H3). For $i \mu \in \sigma(J L) \cap i \mathbf{R} \backslash\{0\}$, there exists a decomposition of $E_{i \mu}$ into closed subspaces $E_{i \mu}=E^{D} \oplus E^{1} \oplus E^{G}$ such that $L$ and $J L$ take the block forms

$$
\langle L \cdot, \cdot\rangle \longleftrightarrow\left(\begin{array}{ccc}
0 & 0 & 0 \\
0 & L_{1} & 0 \\
0 & 0 & L_{G}
\end{array}\right), \quad J L \longleftrightarrow\left(\begin{array}{ccc}
A_{D} & A_{D 1} & A_{D G} \\
0 & i \mu & 0 \\
0 & 0 & A_{G}
\end{array}\right) .
$$

For $\mu=0$, there exists a decomposition $E_{0}=\operatorname{ker} L \oplus E^{D} \oplus E^{1} \oplus E^{G}$ such that $L$ and $J L$ take the block form

$$
\langle L \cdot, \cdot\rangle \longleftrightarrow\left(\begin{array}{cccc}
0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 \\
0 & 0 & L_{1} & 0 \\
0 & 0 & 0 & L_{G}
\end{array}\right), \quad J L \longleftrightarrow\left(\begin{array}{cccc}
0 & A_{0 D} & A_{01} & A_{0 G} \\
0 & A_{D} & A_{D 1} & A_{D G} \\
0 & 0 & 0 & 0 \\
0 & 0 & 0 & A_{G}
\end{array}\right) .
$$

In both cases, all blocks are bounded operators, $L_{1}$ and $L_{G}$ are non-degenerate, $\sigma\left(A_{G}\right)=\sigma\left(A_{D}\right)=\{i \mu\}$, and

$$
\operatorname{dim} E^{G} \leq 3\left(k^{\leq 0}(i \mu)-\operatorname{dim} E^{D}-n^{-}\left(\left.L\right|_{E^{1}}\right)\right), \quad \operatorname{dim} E^{1} \leq \infty
$$

Moreover, $\operatorname{ker}\left(A_{G}-i \mu\right) \subset\left(A_{G}-i \mu\right) E^{G}$, namely, the Jordan canonical form of $J L$ on $E^{G}$ has non-trivial blocks only. Let $1<k_{1}<\cdots<k_{j_{0}}$ be the dimensions of Jordan blocks of $A_{G}$ in $E^{G}$. Suppose there are $l_{j}$ Jordan blocks of size $k_{j} \times k_{j}$. For each $j=1, \ldots, j_{0}$, there exist linearly independent vectors

$$
\begin{equation*}
\left\{u_{p, q}^{(j)} \mid p=1, \ldots, l_{j}, q=1, \ldots, k_{j}\right\} \subset E^{G} \tag{2.14}
\end{equation*}
$$

such that
(1) $\forall 1 \leq p \leq l_{j}$,

$$
\left\{u_{p, q}^{(j)}=(J L-i \mu)^{q-1} u_{p, 1}^{(j)}, q=1, \ldots, k_{j}\right\}
$$

form a Jordan chain of length $k_{j}$. More explicitly,

$$
\left.\begin{array}{rl}
\text { on span }\left\{u_{1,1}^{(j)}, \ldots, u_{1, k_{j}}^{(j)}, \ldots, u_{l_{j}, 1}^{(j)}, \ldots, u_{l_{j}, k_{j}}^{(j)}\right\}: \\
A_{G} \longleftrightarrow\left(\begin{array}{cccccccccc}
i \mu & 0 & \cdots & 0 & 0 & \cdots & 0 & 0 & \cdots & 0 \\
1 & i \mu & \cdots & 0 & 0 & \cdots & 0 & 0 & \cdots & 0 \\
\\
0 & 0 & \cdots & & & 1 & i \mu & \cdots & 0 & 0 \\
\cdots & \cdots & 0 & 0 \\
0 & 0 & \cdots & 0 & 0 & \cdots & i \mu & 0 & \cdots & 0 \\
0 & 0 & \cdots & 0 & 0 & \cdots & 1 & i \mu & \cdots & 0 \\
& & & & & & & \cdots & & \\
0 & 0 & \cdots & 0 & 0 & \cdots & 0 & 0 & \cdots & 1
\end{array}\right) i \mu
\end{array}\right) .
$$

The above count for all Jordan blocks of $A_{G}$ of size $k_{j}$.
(2) $\left\langle L u_{p, q}^{(j)}, u_{p^{\prime}, q^{\prime}}^{\left(j^{\prime}\right)}\right\rangle=0$ if $p \neq p^{\prime}$ or $j \neq j^{\prime}$.
(3) $\forall 1 \leq p \leq l_{j}$, the $k_{j} \times k_{j}$ representation matrix of $L$ on a chain (2.14) is

$$
\left(\left\langle L u_{p, q}^{(j)}, u_{p, r}^{(j)}\right\rangle\right)_{q, r}=\left(\begin{array}{ccccc}
0 & 0 & \cdots & 0 & a_{p, 1}^{(j)}  \tag{2.15}\\
0 & 0 & \cdots & a_{p, 2}^{(j)} & 0 \\
\cdots & & & & \\
a_{p, k_{j}}^{(j)} & 0 & \cdots & 0 & 0
\end{array}\right)
$$

where the entries satisfy

$$
a_{p, q^{\prime}}^{(j)}=(-1)^{q^{\prime}-q} a_{p, q}^{(j)} \neq 0, \quad a_{p, k_{j}+1-q}^{(j)}=\overline{a_{p, q}^{(j)}}
$$

and thus the above matrix is non-degenerate.
(4) If $k_{j}$ is odd, then $a_{p, \frac{1}{2}\left(k_{j}+1\right)}^{(j)}= \pm 1$ and the $k_{j}$-th Krein signature of $i \mu$ defined by

$$
n_{k_{j}}^{-}(i \mu)=\sum_{p=1}^{l_{j}} \min \left\{0, a_{p, \frac{1}{2}\left(k_{j}+1\right)}^{(j)}\right\}
$$

is independent of the choices of such bases $\left\{u_{p, q}^{(j)}\right\}$.

REmARK 2.11. Since $\langle L u, u\rangle$ is symmetric (Hermitian after the complexification), we can normalize the above $a_{p, q}^{(j)}$ such that $a_{p, q}^{(j)}= \pm 1$ if $k_{j}$ is odd and $a_{p, q}^{(j)}= \pm i$ if $k_{j}$ is even. In particular, when $\mu=0$, since the generalized eigenspace is spanned by real functions in $X$, it follows that the Jordan chains in $E^{G} \subset E_{0}$ are all of odd length.

In the splitting of $E_{i \mu}$, we note that only $E^{1}$ may be infinite dimensional, where $L$ is positive except in finitely many directions. If $\langle L \cdot, \cdot\rangle$ is non-degenerate on $E_{i \mu}$, the subspace $E^{D}$ may be eliminated and many of our results can be improved. However, this degeneracy indeed is possible. See such an example in Section 8.4. On the positive side, in that section, we also prove the following proposition on the non-degeneracy of $\left.L\right|_{E_{i \mu}}$ for isolated eigenvalues $i \mu$. In particular, the isolation assumption for $i \mu \in \sigma(J L) \cap i \mathbf{R}$ usually holds if the problem comes from PDEs defined on bounded or periodic domains.

Proposition 2.3. If $i \mu \in \sigma(J L) \cap i \mathbf{R}$ is isolated in $\sigma(J L)$, then
(i) $i \mu$ is an eigenvalue, i.e. $E_{i \mu} \neq\{0\}$, and $\langle L \cdot, \cdot\rangle$ is non-degenerate on $E_{i \mu} /$ $\left(\operatorname{ker} L \cap E_{i \mu}\right)$.
(ii) there exists a closed subspace $E_{\#} \subset X$ invariant under $J L$ such that $X=$ $E_{i \mu} \oplus E_{\#}$ and $\langle L u, v\rangle=0$ for all $u \in E_{i \mu}$ and $v \in E_{\#}$.
(iii) $\sigma\left(\left.(J L)\right|_{E_{\#}}\right)=\sigma(J L) \backslash\{i \mu\}$.

In the case of an isolated spectral point $i \mu$, one may define the invariant eigenspaces and its complement eigenspace via contour integral in operator calculus. Usually it is not guaranteed that such $i \mu$ is an eigenvalue and its eigenspace coincides with $E_{i \mu}$. This proposition implies that, under assumptions (H1-3), this is exactly the case and $\langle L \cdot, \cdot\rangle$ is non-degenerate on $E_{i \mu}$. As a corollary, we prove

Proposition 2.4. In addition to (H1-3), we assume
$(\mathbf{H} 4)\langle L \cdot, \cdot\rangle$ is non-degenerate on $E_{\lambda}$ for any non-isolated $\lambda \in \sigma(J L) \cap i \mathbf{R} \backslash\{0\}$ and also on $E_{0} / \operatorname{ker} L$ if $0 \in \sigma(J L)$ is not isolated,
then there exist closed subspaces $N$ and $M$, which are L-orthogonal, such that $N \oplus \operatorname{ker} L$ and $M \oplus \operatorname{ker} L$ are invariant under $J L, X=N \oplus M \oplus \operatorname{ker} L, \operatorname{dim} N<\infty$, and $L \geq \delta$ on $M$ for some $\delta>0$.

In particular, if eigenvalues of $J L$ are isolated, then by Proposition $2.3,(\mathbf{H} 4)$ is automatically satisfied and Proposition 2.4 holds. If we further assume $\operatorname{ker} L=\{0\}$, then $X=N \oplus M$ and both $N$ and $M$ are invariant under $J L$. Proposition 2.4 can be used to construct invariant decompositions for $L$-self-adjoint operators. The next proposition gives a generalization of Theorem A. 1 in [23], which was proved for a compact $L$-self-adjoint operator $A$ with $\operatorname{ker} A=\{0\}$. Such decomposition was used to study the damping of internal waves in a stably stratified fluid ([23]).

Proposition 2.5. Let $X$ be a complex Hilbert space along with a Hermitian symmetric quadratic form $B(u, v)=\langle L v, u\rangle$ defined by an (anti-linear) operator $L: X \rightarrow X^{*}$ satisfying (H2) with $\operatorname{ker} L=\{0\}$. Let $A: X \rightarrow X$ be a $L$-self-adjoint complex linear operator (i.e. $\langle L A u, v\rangle=\langle L u, A v\rangle$ ) such that nonzero eigenvalues of $A$ are isolated. If $\left.L\right|_{\operatorname{ker} A}$ is non-degenerate, then there exists a decomposition $X=N \oplus M$ such that $N$ and $M$ are L-orthogonal and invariant under $A, \operatorname{dim} N<$ $\infty$ and $\left.L\right|_{M}$ is uniformly positive.

We will extend the notion of the Krein signature to eigenvalues $i \mu$ for which $\langle L \cdot, \cdot\rangle$ on $E_{i \mu}$ is non-degenerate, and give more detailed descriptions of $k_{i}^{-}$and $k_{0}^{-}$. As commented above, the non-degeneracy assumption means $E^{D}$ is eliminated in $E_{i \mu}$. For such $\mu$, define

$$
E_{i \mu, 0}=\left\{v \in \operatorname{ker}(J L-i \mu) \mid\left\langle L v, u_{p, q}^{(j)}\right\rangle=0, \forall 1 \leq j \leq j_{0}, 1 \leq p \leq l_{j}, 1 \leq q \leq k_{j}\right\}
$$

which is the complementary subspace of $R(J L-i \mu) \cap \operatorname{ker}(J L-i \mu)$ inside $\operatorname{ker}(J L-$ $i \mu)$. It corresponds to the diagonalized part of $\left.J L\right|_{E_{i \mu}}$.

Definition 2.1. For $\mu \geq 0$ such that $\langle L \cdot, \cdot\rangle$ is non-degenerate on $E_{i \mu}$, define the first Krein signature

$$
n_{1}^{-}(i \mu)=n^{-}\left(\left.L\right|_{E_{i \mu, 0}}\right)
$$

and $k_{j}$-th Krein signatures as $n_{k_{j}}^{-}(i \mu)$ given in Proposition 2.2, for odd $k_{j}=2 m-$ $1 \geq 1$.

REmark 2.12. The Krein signature $n_{k_{j}}^{-}(i \mu)$, for odd $k_{j}=2 m-1 \geq 1$, does not have to be defined as in Proposition 2.2 using the above special bases. In fact, for any $j$, let $\left\{v_{p, q}^{(j)}\right\}$ be an arbitrary complete set of Jordan chains of length $k_{j}$. Define the $l_{j} \times l_{j}$ matrix $\tilde{M}_{j}=\left(\left\langle L v_{p_{1}, m}^{(j)}, v_{p_{2}, m}^{(j)}\right\rangle\right), 1 \leq p_{1}, p_{2} \leq l_{j}$. Then $n_{k_{j}}^{-}(i \mu)=n^{-}\left(\tilde{M}_{j}\right)$, the negative index (Morse index) of $\tilde{M}_{j}$.

REMARK 2.13. The signatures $n_{k_{j}}^{-}(\mu)$ may also be defined in an intrinsic way independent of bases. See Definition 4.1 and equation (4.2).

According to Proposition 2.2, the 2-dim subspace $\operatorname{span}\left\{u_{p, q}^{(j)}, u_{p, k_{j}+1-q}^{(j)}\right\}$ and 1dim subspace $\operatorname{span}\left\{u_{p, \frac{1}{2}\left(k_{j}+1\right)}^{(j)}\right\}$ for odd $k_{j}$ are $L$-orthogonal to each other. With respect to the basis $\left\{u_{p, q}^{(j)}, u_{p, k_{j}+1-p}^{(j)}\right\}$ there, $L$ takes the form of the Hermitian symmetric matrix $\left(\begin{array}{cc}0 & a \\ \bar{a} & 0\end{array}\right)$ with $a \neq 0$, whose Morse index is clearly 1. Therefore, we obtain the following formula for $k_{i}^{-}$.

Proposition 2.6. In addition to (H1)-(H3), assume $i \mu \in \sigma(J L) \cap i \mathbf{R}$ satisfies that $\langle L \cdot, \cdot\rangle$ is non-degenerate on $E_{i \mu}$. Then we have

$$
\begin{equation*}
k^{\leq 0}(i \mu)=k^{-}(i \mu)=\sum_{k_{j} \text { even }} \frac{l_{j} k_{j}}{2}+\sum_{k_{j} \text { odd }}\left[\frac{l_{j}\left(k_{j}-1\right)}{2}+n_{k_{j}}^{-}(i \mu)\right] \tag{2.16}
\end{equation*}
$$

As Hamiltonian systems often possess additional symmetries which generate nontrivial ker $L, k_{0}^{\leq 0}$ deserves some more discussion if $\operatorname{ker} L \neq\{0\}$. The following propositions are proved in Section 8.3, based on a decomposition of the subspace $E_{0}$. Recall that for any subspace $S \subset X, L$ also induces a quadratic form $\langle L \cdot, \cdot\rangle$ on the quotient space $S /(S \cap \operatorname{ker} L)$.

Proposition 2.7. Assume (H1)-(H3), then $(J L)^{-1}(\operatorname{ker} L)$ is a closed subspace. Furthermore, let

$$
n_{0}=n^{\leq 0}\left(\left.\langle L \cdot, \cdot\rangle\right|_{(J L)^{-1}(\operatorname{ker} L) / \operatorname{ker} L}\right) .
$$

Then
(i) $k_{0}^{\leq 0} \geq n_{0}$.
(ii) If $\langle L \cdot, \cdot\rangle$ is non-degenerate on $(J L)^{-1}(\operatorname{ker} L) / \operatorname{ker} L$, then

$$
k_{0}^{\leq 0}=n_{0}=n^{-}\left(\left.\langle L \cdot, \cdot\rangle\right|_{(J L)^{-1}(\operatorname{ker} L) / \operatorname{ker} L}\right) .
$$

Remark 2.14. Practically, in order to compute $n_{0}$ in the above proposition, let $S \subset(J L)^{-1}(\operatorname{ker} L)$ be a closed subspace such that

$$
\begin{equation*}
(J L)^{-1}(\operatorname{ker} L)=\operatorname{ker} L \oplus S \tag{2.17}
\end{equation*}
$$

then $n_{0}=n^{\leq 0}\left(\left.L\right|_{S}\right)$. Often $S$ can be taken as $(\operatorname{ker} L)^{\perp} \cap(J L)^{-1}(\operatorname{ker} L)$.
It is worth comparing the above results with some classical results (e.g. [29, 30]). Consider a nonlinear Hamiltonian equation

$$
\begin{equation*}
\partial_{t} u=J D H(u) \tag{2.18}
\end{equation*}
$$

which has an additional conserved quantity $P(u)$ (often the momentum, mass etc.) due to some symmetry. Assume that for $c$ in a neighborhood of $c_{0}$, there exists $u_{c}$ such that $D H\left(u_{c}\right)-c D P\left(u_{c}\right)=0$, which gives a relative equilibrium of (2.18) such as traveling waves, standing waves, etc. The linearized equation of (2.18) in some reference frame at $u_{c_{0}}$ takes the form of (2.1) with $L=D^{2} H\left(u_{c_{0}}\right)-c_{0} D^{2} P\left(u_{c_{0}}\right)$. It can be verified that $J D P\left(u_{c_{0}}\right) \in \operatorname{ker} L$ and $\left.L \partial_{c} u_{c}\right|_{c=c_{0}}=D P\left(u_{c_{0}}\right)$. In the case where ker $L=\operatorname{span}\left\{J D P\left(u_{c_{0}}\right)\right\}$ and $J$ is one to one (not necessarily with bounded $J^{-1}$ as assumed in $[29,30]$ ), we have

$$
(J L)^{-1}(\operatorname{ker} L)=\operatorname{span}\left\{J D P\left(u_{c_{0}}\right),\left.\partial_{c} u_{c}\right|_{c=c_{0}}\right\}
$$

when $\left.\frac{d}{d c} P\left(u_{c}\right)\right|_{c=c_{0}} \neq 0$ and

$$
n_{0}=\left\{\begin{array}{ll}
0 & \text { if }\left.\frac{d}{d_{c}} P\left(u_{c}\right)\right|_{c=c_{0}}<0 \\
1 & \text { if }\left.\frac{d}{d c} P\left(u_{c}\right)\right|_{c=c_{0}}>0
\end{array} .\right.
$$

If we further assume $n^{-}(L)=1$, then the combination of Proposition 2.7 and Theorem 2.3 implies the result in [29] that equation (2.1) is stable if $\left.\frac{d}{d c} P\left(u_{c}\right)\right|_{c=c_{0}} \leq$ 0 and unstable if $\left.\frac{d}{d c} P\left(u_{c}\right)\right|_{c=c_{0}}>0$.

In the following special cases, $k_{0}^{\leq 0}$ as well as $n_{0}$ can be better estimated, which is often useful in applications.

Lemma 2.1. Assume (H1)-(H3). we have
(i) $\langle L u, v\rangle=0, \quad \forall u \in \operatorname{ker}(J L), v \in \overline{R(J)}$.
(ii) $\langle L u, u\rangle$ is non-degenerate on $\operatorname{ker}(J L) / \operatorname{ker} L$ if and only if it is non-degenerate on $\overline{R(J)} /(\operatorname{ker} L \cap \overline{R(J)})$.

While the statement of the lemma and the following proposition in the language of quotient spaces make them independent of choices of subspaces transversal to ker $L$, practically it might be easier to work with subspaces. The following is an equivalent restatement of Lemma 2.1 using subspaces. Actually the proof in Section 8.3 will be carried out by using subspaces.

Corollary 2.3. Let $S_{1}, S^{\#} \subset X$ be closed subspaces such that

$$
\begin{equation*}
\operatorname{ker}(J L)=\operatorname{ker} L \oplus S_{1}, \quad \overline{R(J)}=(\overline{R(J)} \cap \operatorname{ker} L) \oplus S^{\#} \tag{2.19}
\end{equation*}
$$

We have that $\langle L \cdot, \cdot\rangle$ is non-degenerate on $S_{1}$ if and only if it is non-degenerate on $S^{\#}$.

Under this non-degeneracy, we have

Proposition 2.8. Assume (H1)-(H3), and that $\langle L u, u\rangle$ is non-degenerate on $\operatorname{ker}(J L) / \operatorname{ker} L$ which is equivalent to $\operatorname{ker}(J L) \cap \overline{R(J)} \subset \operatorname{ker} L$, then
(i) $X=\operatorname{ker}(J L)+\overline{R(J)}$ and

$$
n^{-}(L)=n^{-}\left(\left.L\right|_{\operatorname{ker}(J L) / \operatorname{ker} L}\right)+n^{-}\left(\left.L\right|_{\overline{R(J)} /(\operatorname{ker} L \cap \overline{R(J)})}\right)
$$

(ii) Let

$$
\tilde{S}=\overline{R(J)} \cap(J L)^{-1}(\operatorname{ker} L)
$$

Then

$$
k_{0}^{\leq 0} \geq n^{-}\left(\left.L\right|_{\operatorname{ker}(J L) / \operatorname{ker} L}\right)+n^{\leq 0}\left(\left.L\right|_{\tilde{S} /(\operatorname{ker} L \cap \tilde{S})}\right)
$$

(iii) If, in addition, $\langle L u, u\rangle$ is non-degenerate on $\tilde{S} /(\operatorname{ker} L \cap \tilde{S})$, then

$$
\begin{align*}
& k_{0}^{\leq 0}=n^{-}\left(\left.L\right|_{\operatorname{ker}(J L) / \operatorname{ker} L}\right)+n^{-}\left(\left.L\right|_{\tilde{S} /(\operatorname{ker} L \cap \tilde{S})}\right) \\
& k_{r}+2 k_{c}+2 k_{i}^{\leq 0}=n^{-}\left(\left.L\right|_{\overline{R(J)} /(\operatorname{ker} L \cap \overline{R(J)})}\right)-n^{-}\left(\left.L\right|_{\tilde{S} /(\operatorname{ker} L \cap \tilde{S})}\right) \tag{2.20}
\end{align*}
$$

We notice that the last equality is only a consequence of the previous two equalities on $n^{-}$and $k_{0}^{\leq 0}$ and the index Theorem 2.3.

In terms of subspaces, equivalently we have
Corollary 2.4. Let $S_{1}, S^{\#} \subset X$ be closed subspaces assumed in Corollary 2.3 and $S_{2} \in X$ be a closed subspace such that

$$
\begin{equation*}
\overline{R(J)} \cap(J L)^{-1}(\operatorname{ker} L)=S_{2} \oplus(\overline{R(J)} \cap \operatorname{ker} L) \tag{2.21}
\end{equation*}
$$

Assume the non-degeneracy of $\langle L u, u\rangle$ on $S_{1}$. Under this condition, we have

$$
X=\operatorname{ker} L \oplus S_{1} \oplus S^{\#}
$$

and this decomposition is orthogonal with respect to the quadratic form $\langle L \cdot, \cdot\rangle$. Moreover, we have

$$
n^{-}(L)=n^{-}\left(L_{S_{1}}\right)+n^{-}\left(\left.L\right|_{S^{\#}}\right) \text { and } k_{0}^{\leq 0} \geq n^{-}\left(\left.L\right|_{S_{1}}\right)+n^{\leq 0}\left(\left.L\right|_{S_{2}}\right)
$$

The additional non-degeneracy assumption of $\langle L u, u\rangle$ on $\tilde{S} /(\operatorname{ker} L \cap \tilde{S})$ is equivalent to its non-degeneracy on $S_{2}$ and it implies

$$
\begin{aligned}
& k_{0}^{-}=n^{-}\left(\left.L\right|_{S_{1}}\right)+n^{-}\left(\left.L\right|_{S_{2}}\right) \\
& k_{r}+2 k_{c}+2 k_{i}^{-}=n^{-}\left(\left.L\right|_{S^{\#}}\right)-n^{-}\left(\left.L\right|_{S_{2}}\right)
\end{aligned}
$$

Very often subspaces $S_{1}, S^{\#}, S_{2}$ can be taken as various intersections with $(\operatorname{ker} L)^{\perp}$.

### 2.5. Structural stability/instability

Our next main result is on the spectral properties of the Hamiltonian operator $J L$ under small bounded perturbations. Consider the perturbed linear Hamiltonian system

$$
\begin{equation*}
u_{t}=J_{\#} L_{\#} u, \quad J_{\#}=J+J_{1}, \quad L_{\#}=L+L_{1}, \quad u \in X \tag{2.22}
\end{equation*}
$$

We assume the perturbations satisfy
(A1) $J$ and $L$ satisfies (H1-2) and the perturbations $J_{1}: X^{*} \rightarrow X$ and $L_{1}$ : $X \rightarrow X^{*}$ are bounded operators with $J_{1}^{*}=-J_{1}$ and $L_{1}^{*}=L_{1}$.
(A2) dim ker $L<\infty$;
(A3) $D(J L) \subset D\left(J L_{1}\right)$.

We note that (A2) implies (H3) for $J L$ by Remark 2.3. From the Closed Graph Theorem, $J L_{1}$ is a bounded operator on the Hilbert space $D(J L)$ equipped with the graph norm

$$
\begin{equation*}
\|u\|_{G}^{2} \triangleq\|u\|^{2}+\|J L u\|^{2}, u \in D(J L) ;\left|J L_{1}\right|_{G} \triangleq \sup _{\|u\|_{G}=1}\left\|J L_{1} u\right\| \tag{2.23}
\end{equation*}
$$

We first point out that assumptions (A1-3) imply (H1-3) for $J_{\#} L_{\#}$ when the perturbations are sufficiently small as assumed in Theorem 2.4 below. See Lemma 9.1. As indicated in assumption (A1) we consider bounded perturbations to both the symplectic structure $J$ and the energy quadratic form $L$, while the Hamiltonian structure is preserved. Assumption (A2) ensures $n^{-}\left(L_{\#}\right)<\infty$ so that the perturbed problem is still in our framework. Assumption (A3) is a regularity assumption which implies that $J_{\#} L_{\#}$ is not more unbounded compared to $J L$. Therefore, the resolvent $\left(\lambda-J_{\#} L_{\#}\right)^{-1}$ is a small perturbation of $(\lambda-J L)^{-1}$ as proved in Lemma 9.2.

Let $E^{u, s, c}$ be the unstable/stable/center subspaces of $J L$, as well as the constants $\lambda_{u}>0$, as given in Theorem 2.2. The next theorem and the following proposition will be proved in Section 9.1.

Theorem 2.4. Assume (A1-3). There exist $C, \epsilon_{0}>0$ depending only on $J$ and $L$ such that, if

$$
\begin{equation*}
\epsilon \triangleq\left|J_{1}\right|+\left|L_{1}\right|+\left|J L_{1}\right|_{G} \leq \epsilon_{0} \tag{2.24}
\end{equation*}
$$

then
(a) There exist bounded operators

$$
\begin{align*}
& S_{\#}^{u}: E^{u} \rightarrow E^{s} \oplus E^{c}, \quad S_{\#}^{s}: E^{s} \rightarrow E^{u} \oplus E^{c}, \quad S_{\#}^{c}: E^{c} \rightarrow E^{s} \oplus E^{u} \\
& \quad \text { such that } \\
& \left|S_{\#}^{u, s, c}\right| \leq C \epsilon, \quad e^{t J_{\#} L_{\#}} E_{\#}^{u, s, c}=E_{\#}^{u, s, c}, \quad \text { where } E_{\#}^{u, s, c}=\operatorname{graph}\left(S_{\#}^{u, s, c}\right) \\
& \text { for all } t \in \mathbf{R} . \text { Moreover, } \\
& \left|e^{t J_{\#} L_{\#}}\right|_{E_{\#}^{s}} \mid \leq C\left(1+t^{\operatorname{dim} E^{s}-1}\right) e^{-\left(\lambda_{u}-C \epsilon\right) t}, \quad \forall t \geq 0  \tag{2.25}\\
& \left|e^{t J_{\#} L_{\#}}\right|_{E_{\#}^{u}} \mid \leq C\left(1+|t|^{\operatorname{dim} E^{u}-1}\right) e^{\left(\lambda_{u}-C \epsilon\right) t}, \quad \forall t \leq 0 \\
& \left.\quad\left|e^{t J_{\#} L_{\#}}\right|\right|_{\# \#} ^{c} \mid \leq C e^{C \epsilon|t|}, \quad \forall t \in \mathbf{R} . \tag{2.26}
\end{align*}
$$

(b) $\left\langle L_{\#} \cdot, \cdot\right\rangle$ vanishes on $E_{\#}^{u, s}$, but is non-degenerate on $E_{\#}^{s} \oplus E_{\#}^{u}$, and

$$
E_{\#}^{c}=\left\{u \mid\left\langle L_{\#} u, v\right\rangle=0, \forall v \in E_{\#}^{u} \oplus E_{\#}^{s}\right\}
$$

(c) If $\langle L \cdot, \cdot\rangle \geq \delta>0$ on $E^{c}$, then there exists $C^{\prime}>0$ depending on $\delta$, J, and $L$ such that $\left.\left|e^{t J_{\#} L_{\#}}\right|\right|_{\# \#} ^{c} \mid \leq C^{\prime}$ for any $t \in \mathbf{R}$.
Due to assumption (A3), the resolvent $\left(\lambda-J_{\#} L_{\#}\right)^{-1}$ is only a small perturbation of $(\lambda-J L)^{-1}$ as proved in Lemma 9.2. Therefore, the existence of the invariant subspaces $E_{\#}^{u, s, c}$ as a small perturbation to $E^{u, s, c}$ follows immediately. Statements (b) and (c) basically result from the Hamiltonian structure and the estimates of $e^{t J_{\#} L_{\#}}$ on $E_{\#}^{u, s}$ are basically due to their finite dimensionality. If $J_{\#} L_{\#}-J L$ had been a bounded operator, estimate (2.26) would follow easily from the standard spectral theory as well. However, since $J: X^{*} \supset D(J) \rightarrow X$ is only assumed to satisfy $J^{*}=-J^{*}$, the term $J L_{1}$ may not be bounded and thus (2.26) does not
follow from the standard spectral theory. Our proof heavily relies on the decomposition given by Theorem 2.1. In fact, the usual resolvent estimate often neglects the Hamiltonian structure of the problem which actually plays an essential role here. Otherwise a counterexample without the Hamiltonian structure is $J=J_{\#}=i$ and $L_{\#}=\partial_{x x}+\epsilon \partial_{x}$ with $X=H^{1}\left(S^{1}, \mathbf{C}\right)$, for which the equation $u_{t}=J_{\#} L_{\#} u$ is not even well-posed in $X$ for $\epsilon \neq 0$.

Another consequence of Lemma 9.2 of the resolvent estimate and Lemma 6.2 is the following structural stability type result.

Proposition 2.9. Suppose closed subsets $\sigma_{1,2} \subset \sigma(J L)$ satisfy
(1) $\sigma(J L)=\sigma_{1} \cup \sigma_{2}, \sigma_{1} \cap \sigma_{2}=\emptyset$, and $\sigma_{2}$ is compact.
(2) For any $\lambda \in \sigma_{1}$ and $0 \neq u \in E_{\lambda}$, it holds $\langle L u, u\rangle>0$.

Then there exist $\alpha, \epsilon_{0}>0$ depending only on $J$ and $L$ such that (2.24) implies

$$
\left\{\lambda \in \sigma\left(J_{\#} L_{\#}\right) \mid d\left(\lambda, \sigma_{2}\right)>\alpha\right\} \subset i \mathbf{R}
$$

From Proposition 6.2, any $\lambda \in \sigma(J L) \backslash i \mathbf{R}$ is an eigenvalue, i.e. $E_{\lambda} \neq\{0\}$, and $\langle L \cdot, \cdot\rangle$ vanishes on $E_{\lambda}$. Therefore, it must hold that $\sigma_{1} \subset i \mathbf{R}$. Even though the second assumption on $\sigma_{1}$ seems weaker than that $\langle L \cdot, \cdot\rangle$ is uniformly positive on its eigenspaces, it along with Theorem 2.3 actually implies the latter. This proposition means that, under small perturbations, unstable eigenvalues can not bifurcate from such $\sigma_{1}$.

In the next we consider the deformation of purely imaginary spectral points of $J L$ under perturbations as they are closely related to generation of linear instability. The next two theorems are proved in Section 9.2. Firstly we prove that if $i \mu \in \sigma(J L)$ and $\langle L \cdot, \cdot\rangle$ has certain definite sign on $E_{i \mu}$, then $\sigma\left(J_{\#} L_{\#}\right)$ would not have nearby unstable eigenvalues.

Theorem 2.5. Assume (A1-3), $i \mu \in \sigma(J L) \cap i \mathbf{R}$, and either a.) there exists $\delta>0$ such that $\langle L u, u\rangle \geq \delta\|u\|^{2}$ for all $u \in E_{i \mu}$ or b.) $i \mu$ is isolated in $\sigma(J L)$ and $\langle L u, u\rangle \leq-\delta\|u\|^{2}$ for all $u \in E_{i \mu}$, then there exist $\alpha, \epsilon_{0}>0$ depending on $J, L, \mu$, and $\delta$ such that, if (2.24) holds, then

$$
\left\{\lambda \in \sigma\left(J_{\#} L_{\#}\right)||\lambda-i \mu| \leq \alpha\} \subset i \mathbf{R}\right.
$$

REmark 2.15. On the one hand, note that in the above theorem, we do not require $i \mu$ being an isolated eigenvalue or even an eigenvalue of $J L$. If $i \mu$ is not an eigenvalue, $E_{i \mu}=\{0\}$ and the sign definiteness assumption is automatically satisfied. On the other hand, if $i \mu$ is an isolated spectral point, then Proposition 2.3 implies that $E_{i \mu}$ is nontrivial and is precisely the eigenspace of $i \mu$. Moreover, from Lemma 3.4 and the sign definiteness of $L$ on $E_{i \mu}$, we have $E_{i \mu}=\operatorname{ker}(J L-i \mu)$.

On the one hand, the above theorem indicates that under Hamiltonian perturbations, hyperbolic (i.e. stable and unstable) eigenvalues can not bifurcate from either a.) any $i \mu \in \sigma(J L)$, whether isolated or not, for which $\langle L \cdot, \cdot\rangle$ is positive on $E_{i \mu}$, or b.) any isolated eigenvalue $i \mu$ where $\langle L \cdot, \cdot\rangle$ has a definite sign on $E_{i \mu}$. Theorem 2.5, as well as Theorem 2.4 can be viewed as robustness or structural stability type results.

On the other hand, as given in the next theorem, the structural stability conditions in Theorem 2.5 are also necessary for an eigenvalue $i \mu \neq 0$. As in many applications parameters mostly appear in the energy operator $L$ instead of the symplectic operator $J$, we will study perturbations only to $L$ for possible bifurcations of unstable eigenvalues near $i \mu$.

ThEOREM 2.6. Assume that $(J, L)$ satisfies (H1-3) and $0 \neq i \mu \in \sigma(J L) \cap i \mathbf{R}$ satisfies
(1) $\langle L \cdot, \cdot\rangle$ is neither positive nor negative definite on $E_{i \mu}$ or
(2) $i \mu$ is non-isolated in $\sigma(J L)$ and there exists $u \in E_{i \mu}$ with $\langle L u, u\rangle \leq 0$,
then for any $\epsilon>0$, there exist a symmetric bounded linear operator $L_{1}: X \rightarrow$ $X^{*}$ such that: $\left|L_{1}\right|<\epsilon$ and there exists $\lambda \in \sigma\left(J\left(L+L_{1}\right)\right)$ with $\operatorname{Re} \lambda>0$ and $|\lambda-i \mu|<C \epsilon$, for some constant $C$ depending only on $\mu, J, L$.

It is easy to see that conditions in Theorem 2.6 are exactly complementary to those in Theorem 2.5 for $i \mu \neq 0$ and thus they give necessary and sufficient conditions on whether unstable eigenvalues can bifurcate from $0 \neq i \mu \in \sigma(J L) \cap i \mathbf{R}$ under Hamiltonian perturbations.

Remark 2.16. In [28], Grillakis proved that an embedded purely imaginary eigenvalue with negative energy of the linearized operator at excited states of a semilinear nonlinear Schrödinger equation is 'structurally unstable' under small perturbations and unstable eigenvalues can be generated. The linearized operator is of the form JL, where

$$
J=\left(\begin{array}{cc}
0 & 1 \\
-1 & 0
\end{array}\right), L=\left(\begin{array}{cc}
-\Delta+V_{1}(x) & 0 \\
0 & -\Delta+V_{2}(x)
\end{array}\right) .
$$

Here, $V_{1}(x), V_{2}(x) \rightarrow \omega>0$ exponentially when $|x| \rightarrow \infty$. Under some assumptions, Theorem 2.4 in [28] implies that that if $i \mu \neq 0$ is an embedded eigenvalue of $J L$ with $\langle L u, u\rangle<0$ for some eigenfunction $u$, then an unstable eigenvalue may bifurcate from i $\mu$ under Hamiltonian perturbations. Similar result was also obtained in [21]. This is a special case of the above theorem. Actually, we can relax the structural instability condition to be that $\langle L \cdot, \cdot\rangle$ is not positive definite on $E_{i \mu}$, including cases of degeneracy of $\left.L\right|_{E_{i \mu}}$ or with Jordan chains.

However, it should be pointed out that it is not clear that the above structural instability may be realized by the linearized equation of the nonlinear Schrödinger equation at a perturbed excited state. It would be interesting to see if one can prove the structural instability in the sense that there is linear instability for nearby excited states.

REMARK 2.17. The case $\mu=0$ is not included in Theorem 2.6 since this may be related to some additional degeneracy of $L$ or $J$. See for example Cases $3 b$ and 3d in Section 9.2. The analysis of possible bifurcations of unstable eigenvalues from $\mu=0$ could be carried out in a similar fashion based on the Propositions 2.2, 2.3, Lemma 9.4, etc., but more carefully. We feel that it might be easier to work on this case directly in concrete applications and thus do not include it in the above theorem.

### 2.6. A theorem where $L$ does not have a positive lower bound on $X_{+}$

Among our global assumptions (H1-3), (H2) requires that the phase space $X$ is decomposed into the direct sum of three subspaces $X=X_{-} \oplus \operatorname{ker} L \oplus X_{+}$, such that the quadratic form $\langle L \cdot, \cdot\rangle$ is uniformly positive/negative on $X_{ \pm}$. This assumption plays a crucial role in the analysis throughout the paper. However, in some Hamiltonian PDEs $L$, which usually appears as the Hessian of the energy functional at a steady state, may not have a positive lower bound on $X_{+}$. One such simple example is $X=H^{1}\left(\mathbf{R}^{n}\right)$ and $L=-\Delta+a(x)$ where $\lim _{|x| \rightarrow \infty} a(x)=0$.

Even if $a>0$ which implies $L>0$, but for any $\delta>0$, there exists $u \in H^{1}$ such that $\langle L u, u\rangle<\delta\|u\|_{H^{1}}^{2}$. A potential resolution to this issue in this specific example is to take a different phase space such as $\dot{H}^{1}$ instead of $H^{1}$. In Chapter 10, we show that this observation may be applied in a rather general setting. As a non-trivial example of this case, the stability of traveling waves of a nonlinear Schrödinger equation in 2-dim with non-vanishing condition at $|x|=\infty$ is considered in Section 11.6.

In this section, let $X$ be a real Hilbert space with the inner product $(\cdot, \cdot)$ and we assume
(B1) $Q_{0}, Q_{1}: X \rightarrow X^{*}$ are bounded positive symmetric linear operators such that

$$
\left\langle\left(Q_{0}+Q_{1}\right) u, v\right\rangle=(u, v), Q_{0,1}^{*}=Q_{0,1},\left\langle Q_{0,1} u, u\right\rangle>0, \forall 0 \neq u, v \in X
$$

(B2) $\mathbb{J}: X \rightarrow X$ is a bounded linear operator satisfying

$$
\mathbb{J}^{-1}=-\mathbb{J}, \quad\left\langle Q_{0} \mathbb{J} u, \mathbb{J} u\right\rangle=\left\langle Q_{0} u, u\right\rangle, \quad \forall u \in X
$$

Let $J=\mathbb{J} Q_{0}^{-1}: X^{*} \supset Q_{0}(X) \rightarrow X$.
(B3) $L: X \rightarrow X^{*}$ is a bounded symmetric linear operator such that $L_{1}=L-Q_{1}$ satisfies

$$
\left|\left\langle L_{1} u, v\right\rangle\right|^{2} \leq c_{0}\left(\left\langle Q_{0} u, u\right\rangle\left\langle Q_{0} v, v\right\rangle+\left\langle Q_{0} u, u\right\rangle\left\langle Q_{1} v, v\right\rangle+\left\langle Q_{1} u, u\right\rangle\left\langle Q_{0} v, v\right\rangle\right) .
$$

(B4) There exist closed subspaces $X_{ \pm} \subset X$ such that

$$
\begin{align*}
& X=X_{-} \oplus \operatorname{ker} L \oplus X_{+}, \quad n^{-}(L) \triangleq \operatorname{dim} X_{-}<\infty  \tag{2.27}\\
& \pm\left\langle L u_{ \pm}, u_{ \pm}\right\rangle>0,\left\langle L u_{+}, u_{-}\right\rangle=0, \forall 0 \neq u_{ \pm} \in X_{ \pm} \tag{2.28}
\end{align*}
$$

(B5) Subspaces $X_{ \pm}$satisfy

$$
\operatorname{ker} i_{X_{+}}^{*}=\left\{f \in X^{*} \mid\langle f, u\rangle=0, \forall u \in X_{+}\right\} \subset Q_{0}(X)=D(J)
$$

where $i_{X_{+}}^{*}: X^{*} \rightarrow X_{+}^{*}$ is the dual operator of the embedding $i_{X_{+}}$.
Obviously the assumption in (B1) that $Q_{1}+Q_{0}$ is the Riesz representation of the inner product can be weakened to that it is the Riesz representation of an equivalent inner product. It is also easy to verify that $J$ is closed and antisymmetric, namely, $J \subset-J^{*}$. Roughly the $L$-orthogonal decomposition of $X$ can be constructed a.) by taking $\operatorname{ker} L \oplus X_{+}$as the $L$-orthogonal complement of a carefully chosen $X_{-}$and then $X_{+}$as any complimentary subspace of ker $L$ there; or b.) from a spectral decomposition of the linear operator on $X$ corresponding to the quadratic form $\langle L \cdot, \cdot\rangle$ through certain inner product. In a typical application as in Section 11.6, $Q_{1}$ is often a uniformly positive elliptic operator of order $2 s, L_{1}$ is a perturbation containing lower order derivatives with variable coefficients, and $Q_{0}$ corresponds to the $L^{2}$ duality. It is convenient to start with $X=H^{s}$ initially. The assumption $n^{-}(L)<\infty$ may come from the construction of the steady state via some variational approach. The lack of a positive lower bound of $L$ restricted to $X_{+} \subset H^{s}$ is often due to the missing control of the $L^{2}$ norm by $\langle L \cdot, \cdot\rangle$. This also forces us to make the slightly stronger assumption (B5) than (H3). In Chapter 10 we prove

Theorem 2.7. There exists a Hilbert space $Y$ such that (a) $X$ is densely embedded into $Y$;
(b) $L$ can be extended to a bounded symmetric linear operator $L_{Y}: Y \rightarrow Y^{*}$;
(c) $\left(Y, L_{Y}, J_{Y}\right)$ satisfy (H1-3), where $J_{Y}: D(J) \cap Y^{*} \rightarrow Y$ is the restriction of $J$.

It is natural to define $Y$ through the completion of $X$ under a norm based on $L$. To prove this theorem, the key is to show $(\mathbf{H 1})$ and $(\mathbf{H} 3)$ are satisfied.

### 2.7. Some Applications to PDEs

We briefly discuss the applications of the general theory to several PDE models in Chapter 11. First, we consider the stability of traveling waves of dispersive wave models of KDV, BBM and good Boussinesq types. These PDE models arise as approximation long wave models for water waves etc. We treat general dispersion symbols including nonlocal ones.

For solitary waves, the linearized equations are written in a Hamiltonian form where the symplectic operators $J$ turn out to be non-invertible unbounded operators. The index formula and the exponential trichotomy estimates are obtained from Theorems 2.2 and 2.3.

For periodic waves, the linearized equations for perturbations of the same period are again written in the Hamiltonian form with $J$ having nontrivial kernels. This brings changes to the index counting formula and stability criteria. In recent years, similar index formula had been studied in various cases. Our results give a unified treatment for general dispersion symbols. For both solitary waves and periodic waves, the linear stability conditions are also shown to imply nonlinear orbital stability. For the unstable cases, the exponential dichotomy can be used to show nonlinear instability and even to further construct local invariant (stable, unstable and center) manifolds near the traveling wave orbit in the energy space. Moreover, when a.) the negative dimension of the linearized energy functional is equal to the unstable dimension of the linearized equation and b.) the kernel of the linearized energy functional is generated exactly by the symmetry group of the system, the orbital stability and local uniqueness on the center manifold could be obtained. These invariant manifolds also give a complete description of dynamics near the orbit of unstable profiles. For more details, we refer to recent papers ([37] [38]) on the construction of invariant manifolds near unstable traveling waves of supercritical KDV equation and 3D Gross-Pitavaeskii equation.

We then consider the linearized problems arisen from the modulational (BenjaminFeir, side-band) instability of period waves. Besides obtaining an index formula for each Floquet-Block problem, we also carry out some perturbation analysis to justify that unstable modes in the long wave limit can only arise from zero eigenvalue of the co-periodic problem. Subsequently we obtain the semigroup estimates for both multi-periodic and localized perturbations, which played an important role on the recent proof ([36]) of nonlinear modulational instability of various dispersive models.

As another application, we consider the eigenvalue problem of the form $L u=$ $\lambda u^{\prime}$, which arises in the stability of traveling waves of generalized Bullough-Dodd equation (11.43). Let $J=\partial_{x}^{-1}$, then it is equivalent to the Hamiltonian form $J L u=\lambda u$. Thus general theorems can be applied to get instability index formula and the stability criterion which generalize the results in [69] by relaxing some restrictions. In particular it implies the linear instability of any traveling wave of generalized Bullough-Dodd equation (11.43), removing the convexity assumption in [69].

Next, we consider stability/instability of steady flows of 2D Euler equation in a bounded domain. For a large class of steady flows, the linearized Euler equation can be written in a Hamiltonian form satisfying (H1)-(H3). Here, the symplectic operator $J$ has an infinite dimensional kernel. The index formula is obtained in terms of a reduced operator related to the projection to ker $L$. By using the perturbation theory in Section 2.5, the structural instability in the case of the presence of embedded eigenvalues is shown. The Hamiltonian structures are also useful in studying the enhanced damping and inviscid damping problems.

Lastly, we study the stability of traveling waves of 2D nonlinear Schrödinger equations with nonzero condition at infinity. When written in the Hamiltonian form $J L$, the quadratic form $\langle L \cdot, \cdot\rangle$ does not have uniform lower bound on the positive subspace $X_{+}$. The strategy used for the 3D case ([53]) does not work in 2D. We use the theory in Section 2.6 to construct a new and larger phase space to recover the uniform positivity of $\langle L \cdot, \cdot\rangle$ on the positive space. Then the theory in Section 2.4 is used to prove the stability criterion in terms of the sign of $d P / d c$, where $P(c)$ is the momentum of a traveling wave of speed $c$. As a somewhat unusual application of the index formula, we prove the positivity of the momentum $P$ for traveling waves (in both 2D and 3D) with general nonlinear terms.

## CHAPTER 3

## Basic properties of Linear Hamiltonian systems

In this chapter, we present a few basic qualitative properties of the linear equation (2.1), including the conservation of energy, some elementary spectral properties, etc. As our problem is set up in a functional analysis theoretical framework, in some cases we have to follow the painful rigor at an orthodox level. To make it less tedious, we only keep those basic results directly related to the dynamics of (2.1) in this chapter, while some more elementary properties of (2.1), including its well-posedness (Proposition 12.1), are left in Chapter 12, the Appendix.

Like any Hamiltonian flow, we have the conservation of energy and the symplectic structure of the flow defined by (2.1).

Lemma 3.1 ([59]). For any solutions $u(t), v(t)$ of (2.1), then we have
(1) $\frac{d}{d t}\langle L u(t), v(t)\rangle=0$;
(2) $R(J)$ is invariant under $e^{t J L}$; and
(3) if $J$ is one-to-one ( $J^{-1}$ not necessarily bounded) and $u(0) \in R(J)$, then $\frac{d}{d t}\left\langle J^{-1} u(t), v(t)\right\rangle=0$.
Proof. Property (1) is clearly true if $u(0), v(0) \in D(J L)$ and then the general case follows immediately from a density argument. To prove (2), we first notice, for $x \in D(J L)$,

$$
e^{t J L} x-x=J \int_{0}^{t} L e^{t^{\prime} J L} x d t^{\prime}
$$

Since $J$ is closed and $D(J L)$ is dense, a density argument implies that $\int_{0}^{t} L e^{t^{\prime} J L} x d t^{\prime} \in$ $D(J)$ for all $x$ and the above equality holds for all $x$ and thus $R(J)$ is invariant under $e^{t J L}$. For (3), first consider $v(0) \in D(J L)$ and the above equality yields

$$
\frac{d}{d t}\left\langle J^{-1} u(t), v(t)\right\rangle=\langle L u(t), v(t)\rangle+\left\langle J^{-1} u(t), J L v(t)\right\rangle=0
$$

where we used the assumption that $J$ is anti-self-adjoint. Again the general case of (3) follows from the density of $D(J L)$.

An immediate consequence of the conservation of the quadratic form $\langle L \cdot, \cdot\rangle$ is on invariant subspaces.

Lemma 3.2. Suppose a subspace $X_{1} \subset X$ is invariant under $e^{t J L}$, i.e. $e^{t J L} X_{1} \subset$ $X_{1}$ for all $t \in \mathbf{R}$, then $e^{t J L} X_{2} \subset X_{2}$ for all $t \in \mathbf{R}$ where the closed subspace $X_{2}=\left\{u \in X \mid\langle L u, v\rangle=0, \forall v \in X_{1}\right\}$.

Proof. For any $u \in X_{2}, v \in X_{1}$, and $t \in \mathbf{R}$, Lemma 3.1 and the invariance of $X_{1}$ imply

$$
\left\langle L e^{t J L} u, v\right\rangle=\left\langle L u, e^{-t J L} v\right\rangle=0
$$

which yields the conclusion.

While in a substantial part of the paper, we shall work with the real Hilbert space $X$ and real operators $J, L$, etc., for considerations where complex eigenvalues are involved, we have to work with their standard complexification. See the Appendix (Chapter 12) for details.

Let $\lambda$ be an eigenvalue of $J L$ (i.e. $\lambda \in \sigma(J L)$ ) and

$$
E_{\lambda}=\left\{u \in X \mid(J L-\lambda I)^{k} u=0, \text { for some integer } k \geq 1\right\}
$$

Then by Lemma 3.1, we have
Lemma 3.3 (Lemma 2 in [59] or Lemma 2.7 in [32]). If $v_{1} \in E_{\lambda_{1}}, v_{2} \in E_{\lambda_{2}}$ and $\lambda_{1}+\bar{\lambda}_{2} \neq 0$, then $\left\langle L v_{1}, v_{2}\right\rangle=0$.

The following lemma will be repeatedly used to analyze the structure of $E_{i \mu}$, $\mu \in \mathbf{R}$.

Lemma 3.4. For any $i \mu \in \sigma(J L)(\mu \in \mathbf{R}), u_{1} \in(J L-i \mu)^{l} X$, and $u_{2} \in$ $\operatorname{ker}(J L-i \mu)^{l}$, then $\left\langle L u_{1}, u_{2}\right\rangle=0$.

Proof. First, we observe that for any $u, v \in X$,

$$
\begin{equation*}
\langle L(J L-i \mu) u, v\rangle=-\langle L u,(J L-i \mu) v\rangle \tag{3.1}
\end{equation*}
$$

Let $v \in X$ such that $(J L-i \mu)^{l} v=u_{1}$, then we have

$$
\left\langle L u_{1}, u_{2}\right\rangle=\left\langle L(J L-i \mu)^{l} v, u_{2}\right\rangle=(-1)^{l}\left\langle L v,(J L-i \mu)^{l} u_{2}\right\rangle=0
$$

The following lemma is a direct consequence of Lemma 3.4 and (3.1).
Lemma 3.5. For any $i \mu \in \sigma(J L) \cap i \mathbf{R}$, it holds

$$
E_{i \mu}=\operatorname{ker}(J L-i \mu)^{2 k^{\leq 0}(i \mu)+1}, \mu \neq 0, \quad \text { and } \quad E_{0}=\operatorname{ker}(J L)^{2 k_{0}^{\leq 0}+2}
$$

REmARK 3.1. As $J L-i \mu$ is a generator of a strongly $C^{0}$ semigroup, $(J L-i \mu)^{m}$ is closed for any $m$ and thus $E_{i \mu}$ is a closed subspace and $\left.J L\right|_{E_{i \mu}}$ is a bounded operator with $\sigma\left(\left.J L\right|_{E_{i \mu}}\right)=\{i \mu\}$.

Proof. We first consider $\mu \neq 0$ and argue by contradiction. Suppose $u \in E_{i \mu}$ such that

$$
(J L-i \mu)^{K} u=0, \quad(J L-i \mu)^{K-1} u \neq 0, \quad K \geq 2 k^{\leq 0}(i \mu)+2
$$

For any $K-1 \geq j_{1}, j_{2} \geq K-k^{\leq 0}(i \mu)-1$, we obtain from Lemma 3.4

$$
\left\langle L(J L-i \mu)^{j_{1}} u,(J L-i \mu)^{j_{2}} u\right\rangle=0
$$

Therefore, the quadratic form $\langle L \cdot, \cdot\rangle$ vanishes on $\operatorname{span}\left\{(J L-i \mu)^{K-1} u,(J L-i \mu)^{K-2}, \ldots,(J L-\right.$ $\left.i \mu)^{K-k^{\leq 0}(i \mu)-1} u\right\}$, whose dimension is $k^{\leq 0}(i \mu)+1$. This contradicts the definition of $k^{\leq 0}(i \mu)$.

To finish the proof, we consider $\mu=0$. Again we argue by contradiction. Suppose $u \in E_{0}$ is such that

$$
(J L)^{K} u=0, \quad(J L)^{K-1} u \neq 0, \quad K \geq 2 k_{0}^{\leq 0}+3
$$

Case 1. $(J L)^{K-1} u \notin \operatorname{ker} L$. In this case, clearly

$$
\operatorname{span}\left\{(J L)^{j} u \mid 0 \leq j \leq K-1\right\} \cap \operatorname{ker} L=\{0\}
$$

Let $\tilde{E}_{0} \subset E_{0}$ be a subspace such that $E_{0}=\tilde{E}_{0} \oplus \operatorname{ker} L$ and $(J L)^{j} u \in \tilde{E}_{0}$ for any $0 \leq j \leq K-1$. Much as in the above, $\langle L \cdot, \cdot\rangle$ vanishes on

$$
Z \triangleq \operatorname{span}\left\{(J L)^{K-1} u,(J L)^{K-2} u, \ldots,(J L)^{K-k_{0}^{\leq 0}-1} u\right\} \subset \tilde{E}_{0}
$$

Since $\operatorname{dim} Z=k_{0}^{\leq 0}+1$, this is a contradiction to the definition of $k_{0}^{\leq 0}$.
Cases 2. $(J L)^{K-1} u \in \operatorname{ker} L \backslash\{0\}$. Clearly,

$$
\operatorname{span}\left\{(J L)^{j} u \mid 0 \leq j \leq K-2\right\} \cap \operatorname{ker} L=\{0\}
$$

Let $\tilde{E}_{0} \subset E_{0}$ be a subspace such that $E_{0}=\tilde{E}_{0} \oplus \operatorname{ker} L$ and $(J L)^{j} u \in \tilde{E}_{0}$ for any $0 \leq j \leq K-2$. Let

$$
Z \triangleq \operatorname{span}\left\{(J L)^{K-2} u,(J L)^{K-3} u, \ldots,(J L)^{K-k_{0}^{\leq 0}-2} u\right\} \subset \tilde{E}_{0}
$$

According to Lemma 3.4, for $K-k_{0}^{\leq 0}-2 \leq j_{1}, j_{2} \leq K-2$ and $j_{1}+j_{2} \geq K$, we have $\left\langle L(J L)^{j_{1}} u,(J L)^{j_{2}} u\right\rangle=0$. If $K-k_{0}^{\leq 0}-2 \leq j_{1}, j_{2} \leq K-2$ and $j_{1}+j_{2}<K$, it must hold $j_{1}=j_{2}=K-k_{0}^{\leq 0}-2$ and $K=2 k_{0}^{\leq 0}+3$. Using (3.1) we obtain

$$
\left\langle L(J L)^{K-k_{0}^{\leq 0}-2} u,(J L)^{K-k_{0}^{\leq 0}-2} u\right\rangle=(-1)^{K-k_{0}^{\leq 0}-2}\left\langle L(J L)^{K-1} u, u\right\rangle=0
$$

where in the last equality we used $(J L)^{K-1} u \in \operatorname{ker} L$. Therefore, $\langle L \cdot, \cdot\rangle$ vanishes on $Z$. Since $\operatorname{dim} Z=k_{0}^{\leq 0}+1$, this is again a contradiction to the definition of $k_{0}^{\leq 0}$. The proof of the lemma is complete.

To end the chapter of basic properties, we prove the following Lemma on the symmetry of $\sigma(J L)$ about both axes.

Lemma 3.6. Assume (H1)-(H3), except for $n^{-}(L)<\infty$. Suppose $\lambda \in \sigma(J L)$, then we have
i) $\pm \lambda, \pm \bar{\lambda} \in \sigma(J L)$.
ii) Suppose $\lambda$ is an eigenvalue of $J L$ and assume in addition $\operatorname{ker} L=\{0\}$ or $\lambda \neq 0$, then $\bar{\lambda}$ is also an eigenvalue of $J L$ and $-\lambda,-\bar{\lambda}$ are eigenvalues of $(J L)^{*}=$ $-L J$. Moreover, for any $k>0$,

$$
\begin{equation*}
\operatorname{ker}(J L-\bar{\lambda})^{k}=\left\{\bar{u} \mid u \in \operatorname{ker}(J L-\lambda)^{k}\right\} \tag{3.2}
\end{equation*}
$$

and

$$
\begin{equation*}
L: \operatorname{ker}(J L-\lambda)^{k} \rightarrow \operatorname{ker}\left((J L)^{*}+\bar{\lambda}\right)^{k}=\operatorname{ker}(L J-\bar{\lambda})^{k} \tag{3.3}
\end{equation*}
$$

is an anti-linear isomorphism.
iii) Suppose $\lambda$ is an isolated eigenvalue of $J L$ with finite algebraic multiplicity, then $-\lambda, \pm \bar{\lambda}$ are also eigenvalues of $J L$ with the same algebraic and geometric multiplicities.

Here the operators $J$ and $L$ are understood as their complexification, and thus are anti-linear mappings satisfying (12.9).

Proof. As i) is trivial if $\lambda=0$, so we assume $\lambda \neq 0$ or $\operatorname{ker} L=\{0\}$.
Due to (12.12) which states that $J L$ is real, (3.2) and $\bar{\lambda} \in \sigma(J L)$ follow immediately. We are left to prove $-\lambda,-\bar{\lambda} \in \sigma(J L)$ and (3.3).

The anti-linearity property (12.9) implies

$$
\begin{equation*}
L(J L-\lambda) u=(L J-\bar{\lambda}) L u=-\left((J L)^{*}+\bar{\lambda}\right) L u, \forall u \in D(J L) \tag{3.4}
\end{equation*}
$$

Therefore, we have that, for any integer $k>0$,

$$
\begin{equation*}
\left((J L)^{*}+\bar{\lambda}\right)^{k} L u=(-1)^{k} L(J L-\lambda)^{k} u, \forall u \in D\left((J L)^{k}\right) \tag{3.5}
\end{equation*}
$$

It follows from (3.5) that $L\left(\operatorname{ker}(J L-\lambda)^{k}\right) \subset \operatorname{ker}\left((J L)^{*}+\bar{\lambda}\right)^{k}$. Under the assumption $\lambda \neq 0$ or $\operatorname{ker} L=\{0\}$, it holds $\operatorname{ker} L \cap E_{\lambda}=\{0\}$ and thus $L$ is one-to-one on $E_{\lambda}$. Therefore, if $\lambda$ is an eigenvalue of $J L$, then $E_{\lambda}$ is nontrivial which implies $L\left(\operatorname{ker}(J L-\lambda)^{k}\right)$, as well as $\operatorname{ker}\left((J L)^{*}+\bar{\lambda}\right)^{k}$, are nontrivial. We obtain that $-\bar{\lambda}$, as well as $-\lambda$, is an eigenvalue of $(J L)^{*}$. Consequently $-\lambda,-\bar{\lambda} \in \sigma(J L)$.

To finish the proof of (3.3), we only need to show

$$
L\left(\operatorname{ker}(J L-\lambda)^{k}\right) \supset \operatorname{ker}\left((J L)^{*}+\bar{\lambda}\right)^{k}
$$

This is obvious from (3.5) if ker $L=\{0\}$. In the case of $\lambda \neq 0$, it is clear ker $\left((J L)^{*}+\right.$ $\bar{\lambda})^{k} \subset R(L)$. Therefore, for any $v \in \operatorname{ker}\left((J L)^{*}+\bar{\lambda}\right)^{k}$, there exists $u_{1} \in X$ such that $v=L u_{1}$. Equation (3.5) again implies

$$
w=(J L-\lambda)^{k} u_{1} \in \operatorname{ker} L
$$

As $\lambda \neq 0$, let $u=u_{1}-(-\lambda)^{-k} w$, then since $(J L-\lambda) w=(-\lambda) w$, we have $v=L u$ and $u \in \operatorname{ker}(J L-\lambda)^{k}$ due to $(J L-\lambda)^{k} w=(-\lambda)^{k} w$. Therefore, $\operatorname{ker}\left((J L)^{*}+\bar{\lambda}\right)^{k} \subset$ $L \operatorname{ker}(J L-\lambda)^{k}$ and thus $L$ is an one to one correspondence (actually an anti-linear isomorphism) from $\operatorname{ker}(J L-\lambda)^{k}$ to $\operatorname{ker}\left((J L)^{*}+\bar{\lambda}\right)^{k}$.

Finally, suppose $\lambda \neq 0$ and $R(J L-\lambda) \neq X$, we will show $R\left((J L)^{*}+\bar{\lambda}\right) \neq X^{*}$ which implies $-\bar{\lambda},-\lambda \in \sigma\left((J L)^{*}\right)=\overline{\sigma(J L)}$ and thus completes the proof of (i). Assume one the contrary $R\left((J L)^{*}+\bar{\lambda}\right)=X^{*}$. Let $\gamma \in D(J) \backslash R(L)$, according to Remark 2.4, there exists $u \in \operatorname{ker} L$ such that $\langle\gamma, u\rangle \neq 0$. One can compute $\left\langle\left((J L)^{*}+\bar{\lambda}\right) \gamma, u\right\rangle=\bar{\lambda}\langle\gamma, u\rangle \neq 0$ and thus $\left((J L)^{*}+\bar{\lambda}\right) \gamma \notin R(L)$. Therefore, if $R\left((J L)^{*}+\bar{\lambda}\right)=X^{*}$, it must hold $\left((J L)^{*}+\bar{\lambda}\right)(R(L))=R(L)$, which is the range of the right side of (3.4). However, since $\lambda \neq 0$, we have $(J L-\lambda)(\operatorname{ker} L)=\operatorname{ker} L$. Along with $R(J L-\lambda) \neq X$, it implies $R(L) \not \subset R(L(J L-\lambda)$ ), which is the range of the left side of (3.4). We obtain a contradiction and thus $R\left((J L)^{*}+\bar{\lambda}\right) \neq X^{*}$.

If $\lambda \in \sigma(J L)$ is isolated and of finite multiplicity, then the same is true for $\bar{\lambda}$. By i) and ii), $-\lambda,-\bar{\lambda} \in \sigma\left((J L)^{*}\right)$ are also isolated and of the same multiplicities, this implies that $-\lambda,-\bar{\lambda} \in \sigma(J L)$ have the same (geometric and algebraic) multiplicities (see [42] P. 184).

Remark 3.2. As in the proof of Lemma 12.3 and Corollary 12.3, the assumption $n^{-}(L)<\infty$ is not required in the above proof. So Lemma 3.6 holds even when $n^{-}(L)=\infty$. On the other hand, this lemma gives the symmetry of $\sigma(J L)$, but not for general eigenvalues, except for purely imaginary eigenvalues or isolated eigenvalues of finite multiplicity. If $\lambda \in \sigma(J L)$ is a nonzero eigenvalue which is non-isolated or of infinite multiplicity, then above lemma implies that $-\lambda,-\bar{\lambda}$ are eigenvalues of $\sigma\left((J L)^{*}\right)$. In general, we can not exclude the possibility that $-\lambda,-\bar{\lambda}$ are not eigenvalues of $J L$. However, when $n^{-}(L)<\infty$, any $\lambda \in \sigma(J L)$ with $\operatorname{Re} \lambda \neq 0$ must be isolated and of finite multiplicity, and the symmetry of eigenvalues and the dimensions of their eigenspaces are given in Corollary 6.1.

## CHAPTER 4

## Finite dimensional Hamiltonian systems

In this chapter, we consider the case where the energy space $X$ of (2.1) is $X=\mathbf{R}^{n}$ which is complexified to $\mathbf{C}^{n}$. The assumptions (H.1-3) become that $J$ is a real anti-symmetric $n \times n$ matrix and $L$ is a real symmetric $n \times n$ matrix. The counting formula (2.13) essentially follows from [59], except for the formula (2.16). We do not need to assume that $J$ is invertible as assumed in [59].

For $\lambda \in \sigma(J L)$, define

$$
I_{\lambda}=E_{\lambda} \oplus E_{-\bar{\lambda}} \text { if } \lambda \notin i \mathbf{R}, \text { and } I_{\lambda}=E_{\lambda} \text { if } \lambda \in i \mathbf{R}
$$

We have $\mathbf{C}^{n}=I_{\lambda_{1}} \oplus \cdots \oplus I_{\lambda_{l}}$, where $\lambda_{j} \in \sigma(J L)$ are all distinct eigenvalues of $J L$ with $\operatorname{Re} \lambda_{j} \geq 0$. By Lemma 3.3, we have

$$
\begin{equation*}
n^{-}(L)=\sum_{j} n^{-}\left(\left.L\right|_{I_{\lambda_{j}}}\right) \tag{4.1}
\end{equation*}
$$

Based on Lemma 3.6 of the symmetry of $\sigma(J L)$, to prove Theorem 2.3 in the finite dimensional case, it suffices to compute $n^{-}\left(\left.L\right|_{I_{\lambda}}\right)$ for any $0 \neq \lambda \in \sigma(J L) \backslash i \mathbf{R}$.

Lemma 4.1 ([59]). Let $\lambda \in \sigma(J L)$. Assume $\operatorname{ker} L=\{0\}$ or $\lambda \neq 0$, then the restriction $\left.\langle L \cdot, \cdot\rangle\right|_{I_{\lambda}}$ is non-degenerate.

Proof. Suppose $\left.\langle L \cdot, \cdot\rangle\right|_{I_{\lambda}}$ is degenerate. Then there exists $0 \neq u \in I_{\lambda}$ such that $\langle L u, v\rangle=0$ for any $v \in I_{\lambda}$. Since $\mathbf{C}^{n}$ is the direct sum of all different $I_{\lambda^{\prime}}$, $\lambda^{\prime} \in \sigma(J L)$, this implies that $\langle L u, v\rangle=0$ for any $v \in \mathbf{C}^{n}$ by Lemma 3.3. So $L u=0$ and thus $0 \neq u \in I_{\lambda} \cap$ ker $L$. It implies that $\lambda=0$ and $\operatorname{ker} L \neq\{0\}$, a contradiction to our assumptions.

Lemma 4.2 ([59] or [32]). If $\operatorname{Re} \lambda>0$ and let $m_{\lambda}$ to be the algebraic multiplicity of $\lambda$. Then $n^{-}\left(\left.L\right|_{I_{\lambda}}\right)=m_{\lambda}$.

Proof. From Lemma 3.3, the quadratic form $\langle L \cdot, \cdot\rangle$ on $I_{\lambda}=E_{\lambda} \oplus E_{-\bar{\lambda}}$ can be represented in the block form $\left(\begin{array}{cc}0 & A \\ A^{*} & 0\end{array}\right)$. Lemma 4.1 implies the non-degeneracy of $A$ and thus the lemma follows.

The counting formula (2.13) in the finite dimensional case follows from these lemmas and (4.1).

In the rest of this section, we carefully analyze $k^{-}(i \mu)=n^{-}\left(\left.L\right|_{E_{i \mu}}\right), \mu \in \mathbf{R}$, and obtain Proposition 2.2 in finite dimensions. Based on Lemma 3.4 and equation (3.1), we first prove

Lemma 4.3. Suppose $i \mu \in \sigma(J L)(\mu \in \mathbf{R})$ and $K>0$ is an integer, then
(1) for $u, v \in \operatorname{ker}(J L-i \mu)^{K}$,

$$
Q_{K}(u, v) \triangleq i^{K-1}\left\langle L(J L-i \mu)^{K-1} u, v\right\rangle
$$

defines a Hermitian form on $\operatorname{ker}(J L-i \mu)^{K}$; and
(2) assume $\operatorname{ker} L=\{0\}$ or $\mu \neq 0$, then

$$
Y_{K} \triangleq\left(\operatorname{ker}(J L-i \mu)^{K} \cap R(J L-i \mu)\right)+\operatorname{ker}(J L-i \mu)^{K-1}=\operatorname{ker} Q_{K}
$$

Proof. That $Q_{K}$ is a Hermitian form on $\operatorname{ker}(J L-i \mu)^{K}$ is an immediate consequence of equation (3.1). Lemma 3.4 also implies $Y_{K} \subset \operatorname{ker} Q_{K}$. We will show $Y_{K}=\operatorname{ker} Q_{K}$ under the additional assumption $\operatorname{ker} L=\{0\}$ or $\mu \neq 0$. Suppose $u \in \operatorname{ker}(J L-i \mu)^{K}$ is such that

$$
Q_{K}(u, v)=i^{K-1}\left\langle L(J L-i \mu)^{K-1} u, v\right\rangle=0, \quad \forall v \in \operatorname{ker}(J L-i \mu)^{K}
$$

By duality, it implies that

$$
\left.L(J L-i \mu)^{K-1} u \in\left((J L-i \mu)^{K}\right)\right)^{*}\left(\mathbf{C}^{n}\right)=(L J-i \mu)^{K}\left(\mathbf{C}^{n}\right)
$$

Therefore, there exists $w \in \mathbf{C}^{n}$ such that

$$
L(J L-i \mu)^{K-1} u=(L J-i \mu)^{K} w
$$

Since $\mu \neq 0$ or $L$ is surjective, the above equation implies $w \in R(L)$ and thus there exists $\tilde{w} \in \mathbf{C}^{n}$ such that $w=L \tilde{w}$. Consequently,

$$
(L J-i \mu)^{K-1} L(u-(J L-i \mu) \tilde{w})=L(J L-i \mu)^{K-1} u-(L J-i \mu)^{K} w=0
$$

which along with Lemma 3.6 implies

$$
L(u-(J L-i \mu) \tilde{w}) \in \operatorname{ker}(L J-i \mu)^{K-1}=L \operatorname{ker}(J L-i \mu)^{K-1}
$$

Therefore, there exists $v \in \operatorname{ker}(J L-i \mu)^{K-1}$ such that

$$
y=u-(J L-i \mu) \tilde{w}-v \in \operatorname{ker} L
$$

If $\mu \neq 0$, let $w_{1}=\tilde{w}+\frac{1}{i \mu} y$. If $\operatorname{ker} L=\{0\}$, we have $y=0$ and let $w_{1}=\tilde{w}$. In both cases, we have

$$
u=v+(J L-i \mu) w_{1}, \quad v \in \operatorname{ker}(J L-i \mu)^{K-1} \subset \operatorname{ker}(J L-i \mu)^{K}
$$

Therefore, $(J L-i \mu) w_{1} \in R(J L-i \mu) \cap \operatorname{ker}(J L-i \mu)^{K}$ and then $u \in Y_{K}$. The proof is complete.

Corollary 4.1. Assume $\operatorname{ker} L=\{0\}$ or $\mu \neq 0$, then $Q_{K}$ induces a nondegenerate Hermitian form on the quotient space $\operatorname{ker}(J L-i \mu)^{K} / Y_{K}$.

If $K$ is odd, for any $u, v \in \operatorname{ker}(J L-i \mu)^{K}$, clearly

$$
\begin{equation*}
Q_{K}(u, v)=\left\langle L(J L-i \mu)^{\frac{K-1}{2}} u,(J L-i \mu)^{\frac{K-1}{2}} v\right\rangle \tag{4.2}
\end{equation*}
$$

Definition 4.1. For odd $K$, define $n_{K}^{-}(i \mu)$ to be the negative index of the quadratic form $Q_{K}$.

The above quotient space $\operatorname{ker}(J L-i \mu)^{K} / Y_{K}$ is closely related to Jordan chains. Suppose a basis of $\mathbf{C}^{n}$ realizes the Jordan canonical form of $J L$, and there are totally $l$ Jordan blocks of size $K \times K$ corresponding to $i \mu$. There must be $l$ Jordan chains of length $K$ in such basis, each of which is generated by some $v \in \operatorname{ker}(J L-i \mu)^{K} / Y_{K}$ as

$$
v,(J L-i \mu) v, \ldots,(J L-i \mu)^{K-1} v
$$

From standard linear algebra, we have the following lemma.

Lemma 4.4. Vectors $v_{1,1}, \ldots, v_{l, 1}$ generate all $l$ Jordan chains of length $K$ in the sense that

$$
v_{j, k}=(J L-i \mu)^{k-1} v_{j, 1}, \quad 1 \leq k \leq K, 1 \leq j \leq l
$$

are in a basis of $\mathbf{C}^{n}$ realizing all l Jordan blocks of size $K$ of $J L$ corresponding to $i \mu \in \sigma(J L)$, if and only if

$$
v_{1,1}+Y_{K}, \ldots, v_{l, 1}+Y_{K}
$$

form a basis of $\operatorname{ker}(J L-i \mu)^{K} / Y_{K}$.
The following lemma would lead to the realization of the Jordan canonical form of $J L$ and skew-diagonalization of $L$ simultaneously.

Lemma 4.5. Assume $\operatorname{ker} L=\{0\}$ or $\mu \neq 0$ where $i \mu \in \sigma(J L) \cap i \mathbf{R}$. Suppose $\operatorname{dim} \operatorname{ker}(J L-i \mu)^{K} / Y_{K}=l>0$ and $Z \subset \operatorname{ker}(J L-i \mu)^{K}$ satisfies

$$
J L(Z)=Z \quad \text { and } Z /\left(Y_{K} \cap Z\right)=\operatorname{ker}(J L-i \mu)^{K} / Y_{K}
$$

then there exist $v_{1}, \ldots, v_{l} \in Z$ such that

$$
\begin{equation*}
\left\langle L(J L-i \mu)^{m} v_{j}, v_{k}\right\rangle= \pm i^{K-1} \delta_{j, k} \delta_{m, K-1}, \quad 0 \leq m \leq K-1 \tag{4.3}
\end{equation*}
$$

Proof. Since $Q_{K}$ induces a non-degenerate Hermitian form on $\operatorname{ker}(J L-i \mu)^{K} /$ $Y_{K}=Z /\left(Z \cap Y_{K}\right)$, there exist $w_{1}, \ldots, w_{l} \in Z$ such that $w_{1}+Y_{K}, \ldots, w_{l}+Y_{K}$ form a basis of $\operatorname{ker}(J L-i \mu)^{K} / Y_{K}$ and diagonalize $Q_{K}$, that is,

$$
\left\langle L(J L-i \mu)^{K-1} w_{j}, w_{k}\right\rangle=(-i)^{K-1} Q_{K}\left(w_{j}, w_{k}\right)= \pm i^{K-1} \delta_{j, k}
$$

Therefore, we have found $w_{1}, \ldots, w_{l}$ satisfying (4.3) for $m=K-1$.
Suppose $1 \leq m_{0}+1 \leq K-1$ and we have found $w_{1}, \ldots, w_{l} \in Z$ satisfying (4.3) for $m \geq m_{0}+1$. Denote

$$
\alpha=K-1-m_{0} \geq 1, Q_{K}\left(w_{j}, w_{k}\right)=b_{j, k}= \pm \delta_{j, k},\left\langle L(J L-i \mu)^{m_{0}} w_{j}, w_{k}\right\rangle=c_{j, k}
$$

In the next step we will construct $v_{1}, \ldots, v_{l}$ satisfying (4.3) for $m \geq m_{0}$ in the form of

$$
v_{j}=w_{j}+\sum_{j^{\prime}=1}^{j} a_{j, j^{\prime}}(J L-i \mu)^{\alpha} w_{j^{\prime}} \in Z
$$

According to $(3.1),\left\langle L(J L-i \mu)^{m} \cdot, \cdot\right\rangle$ is Hermitian or anti-Hermitian. Without loss of generality, we may consider only $j \leq k$ in (4.3). Compute using (3.1)

$$
\begin{align*}
\langle L(J L & \left.-i \mu)^{m} v_{j}, v_{k}\right\rangle=(-1)^{\alpha} \sum_{k^{\prime}=1}^{k} \overline{a_{k, k^{\prime}}}\left\langle L(J L-i \mu)^{\alpha+m} w_{j}, w_{k^{\prime}}\right\rangle \\
& +\left\langle L(J L-i \mu)^{m} w_{j}, w_{k}\right\rangle+\sum_{j^{\prime}=1}^{j} a_{j, j^{\prime}}\left\langle L(J L-i \mu)^{\alpha+m} w_{j^{\prime}}, w_{k}\right\rangle  \tag{4.4}\\
& +(-1)^{\alpha} \sum_{j^{\prime}=1}^{j} \sum_{k^{\prime}=1}^{k} a_{j, j^{\prime}} \overline{a_{k, k^{\prime}}}\left\langle L(J L-i \mu)^{2 \alpha+m} w_{j^{\prime}}, w_{k^{\prime}}\right\rangle
\end{align*}
$$

If $m+\alpha=K-1-m_{0}+m \geq K$, the induction assumption and the above equation imply

$$
\left\langle L(J L-i \mu)^{m} v_{j}, v_{k}\right\rangle=\left\langle L(J L-i \mu)^{m} w_{j}, w_{k}\right\rangle= \pm i^{K-1} \delta_{j, k} \delta_{m, K-1}
$$

and thus (4.3) for $m \geq m_{0}+1$ holds for these $v_{1}, \ldots, v_{l}$ with any choices of $a_{j, j^{\prime}}$. For $m=m_{0}$, i.e. $m+\alpha=K-1$, if $j<k$, (4.4) implies

$$
\left\langle L(J L-i \mu)^{m_{0}} v_{j}, v_{k}\right\rangle=c_{j, k}+(-1)^{\alpha}(-i)^{K-1} \overline{a_{k, j}} b_{j, j} .
$$

Noticing $b_{j, j}= \pm 1$ and letting

$$
a_{k, j}=(-1)^{\alpha+1}(-i)^{K-1} \overline{c_{j, k}} b_{j, j}, j<k
$$

then we have

$$
\left\langle L(J L-i \mu)^{m_{0}} v_{j}, v_{k}\right\rangle=0, \quad j<k
$$

If $j=k$,

$$
\left.\left\langle L(J L-i \mu)^{m_{0}} v_{k}, v_{k}\right\rangle=c_{k, k}+(-i)^{K-1} b_{k, k}\left(a_{k, k}+(-1)^{\alpha} \overline{a_{k, k}}\right)\right)
$$

Let

$$
a_{k, k}=-\frac{1}{2} i^{K-1} b_{k, k} c_{k, k}=-\frac{1}{2} i^{\alpha} b_{k, k} i^{m_{0}} c_{k, k}
$$

Since (3.1) implies $c_{k, j}=(-1)^{m_{0}} \overline{c_{j, k}}$, we have $i^{m_{0}} c_{k, k} \in \mathbf{R}$ and thus $i^{\alpha} a_{k, k} \in \mathbf{R}$ which makes it easy to verify

$$
\left\langle L(J L-i \mu)^{m_{0}} v_{k}, v_{k}\right\rangle=0 .
$$

Therefore, $v_{1}, \ldots, v_{l} \in Z$ satisfy (4.3) for all $m \geq m_{0}$ and the lemma follows from the induction.

We are in a position to prove Proposition 2.2 in finite dimensions.

Proof of Proposition 2.2 assuming $\operatorname{dim} X<\infty$ and $\operatorname{ker} L=\{0\}$ : Let $E^{D}=\{0\}$, then $1<k_{1}<\cdots<k_{j_{0}}$ are the dimensions of nontrivial Jordan blocks in $E_{i \mu}, \mu \in \mathbf{R}$, and there are $l_{j}>0$ Jordan blocks of size $k_{j}$. For each $1 \leq j \leq j_{0}$, we will find linearly independent

$$
\left\{u_{p, q}^{(j)} \mid p=1, \ldots, l_{j}, q=1, \ldots, k_{j}\right\} \subset E_{i \mu}
$$

which form all Jordan chains of length $k_{j}$ and satisfy the desired properties. The construction is by induction on $j$.

For $j=j_{0}$, applying Lemma 4.5 to $Z=\operatorname{ker}(J L-i \mu)^{k_{j}}=E_{i \mu}$, where $\operatorname{dim} \operatorname{ker}(J L-i \mu)^{k_{j_{0}}} / Y_{k_{j_{0}}}=l_{j_{0}}$ according to Lemma 4.4, there exist $u_{1,1}^{\left(j_{0}\right)}, \ldots, u_{l_{j_{0}, 1}}^{\left(j_{0}\right)}$ such that

$$
\begin{equation*}
\left\langle L(J L-i \mu)^{m} u_{p_{1}, 1}^{\left(j_{0}\right)}, u_{p_{2}, 1}^{\left(j_{0}\right)}\right\rangle= \pm i^{k_{j_{0}}-1} \delta_{p_{1}, p_{2}} \delta_{m, k_{j_{0}}-1}, \quad 0 \leq m \leq j_{0}-1 \tag{4.5}
\end{equation*}
$$

In particular we have $Q_{K}\left(u_{p_{1}, 1}^{\left(j_{0}\right)}, u_{p_{2}, 1}^{\left(j_{0}\right)}\right)= \pm \delta_{p_{1}, p_{2}}$. Lemma 4.3 and Corollary 4.1 imply that $u_{1,1}^{\left(j_{0}\right)}+Y_{k_{j_{0}}}, \ldots, u_{l_{j_{0}}, 1}^{\left(j_{0}\right)}+Y_{k_{j_{0}}}$ form a basis of $\operatorname{ker}(J L-i \mu)^{k_{j_{0}}} / Y_{k_{j_{0}}}$. From Lemma 4.4, we obtain that

$$
u_{p, q}^{\left(j_{0}\right)}=(J L-i \mu)^{q-1} u_{p, 1}^{\left(j_{0}\right)}, \quad q=1, \cdots, k_{j_{0}}, p=1, \ldots, l_{j_{0}}
$$

form $l_{j_{0}}$ Jordan chains realizing all Jordan blocks of size $k_{j_{0}}$ of $J L$ corresponding to $i \mu \in \sigma(J L)$. Moreover, equation (4.5) implies

$$
\left\langle L u_{p_{1}, q_{1}}^{\left(j_{0}\right)}, u_{p_{2}, q_{2}}^{\left(j_{0}\right)}\right\rangle= \pm i^{k_{j_{0}}-1} \delta_{p_{1}, p_{2}} \delta_{q_{1}+q_{2}, k_{j_{0}}+1}
$$

Suppose $0 \leq j_{*}<j_{0}$ and we have constructed linearly independent $u_{p, q}^{(j)}$ for all $j_{*}<j \leq j_{0}, 1 \leq p \leq l_{j}, 1 \leq q \leq k_{j}$ satisfying

$$
\begin{equation*}
u_{p, q}^{(j)}=(J L-i \mu)^{q-1} u_{p, 1}^{(j)}, \quad\left\langle L u_{p_{1}, q_{1}}^{(j)}, u_{p_{2}, q_{2}}^{(j)}\right\rangle= \pm i^{k_{j}-1} \delta_{p_{1}, p_{2}} \delta_{q_{1}+q_{2}, k_{j}+1} \tag{4.6}
\end{equation*}
$$

Clearly,

$$
Z_{1}=\operatorname{span}\left\{u_{p, q}^{(j)} \mid j_{*}<j \leq j_{0}, 1 \leq p \leq l_{j}, 1 \leq q \leq k_{j}\right\} \subset E_{i \mu}
$$

is a subspace invariant under $J L$. Moreover, vectors $\left\{u_{p, q}^{(j)}\right\}$ form a basis of $Z_{1}$ realizing the Jordan canonical form of $J L$ on $Z_{1}$ consisting of all those Jordan blocks of $J L$ corresponding to $i \mu$ of size greater than $k_{j_{*}}$. According to (4.6), the quadratic form $\langle L \cdot, \cdot\rangle$ is non-degenerate on $Z_{1}$. In the next step we will construct $u_{p, q}^{\left(j_{*}\right)}$ for $1 \leq p \leq l_{j_{*}}$ and $1 \leq q \leq k_{j_{*}}$. Let

$$
Z=\left\{u \in E_{i \mu} \mid\langle L u, v\rangle=0, \forall v \in Z_{1}\right\} .
$$

Due to the non-degeneracy of $\langle L \cdot, \cdot\rangle$ on both $Z_{1}$ and $I_{i \mu}=E_{i \mu}$ (Lemma 4.1), we have $E_{i \mu}=Z_{1} \oplus Z$. For any $u \in Z$ and $v \in Z_{1}$, due to the symmetry of $L$ and $J$, we have

$$
\langle L J L u, v\rangle=-\langle L u, J L v\rangle=0, \quad \text { as } J L v \in Z_{1}
$$

which implies $J L(Z) \subset Z$. Since the Jordan canonical form of $J L$ on $Z_{1}$ includes all Jordan blocks of $J L$ on $E_{i \mu}$ of size greater than $k_{j_{*}}$, the Jordan canonical form of $J L$ on $Z$ must be those Jordan blocks of $J L$ on $E_{i \mu}$ of size no greater than $k_{j_{*}}$. Therefore, $Z \subset \operatorname{ker}(J L-i \mu)^{k_{j_{*}}}$ and then Lemma 4.4 implies $Z /\left(Z \cap Y_{k_{j_{*}}}\right)=$ $\operatorname{ker}(J L-i \mu)^{k_{j_{*}}} / Y_{k_{j_{*}}}$. Lemma 4.5 provides vectors $u_{1,1}^{\left(j_{*}\right)}, \ldots, u_{l_{j_{*}}, 1}^{\left(j_{*}\right)} \in Z$. It is easy to verify that $u_{p, q}^{(j)}, j_{*} \leq j \leq j_{0}$, satisfy the induction assumption for $j_{*} \leq j \leq j_{0}$. Therefore, by induction, we find all $u_{p, q}^{(j)}$ satisfying (4.6) and realizing all Jordan blocks of $J L$ on $E_{i \mu}$ of size greater than 1 . It is straightforward to verify all the properties in Proposition 2.2. In particular, Lemma 4.4 and equation (4.2) imply that the Krein signature defined in Proposition 2.2 and Remark 2.12 coincides with the one in the above Definition 4.1 in terms of $Q_{K}$. Therefore, it is independent of the choice of the basis (Jordan chains) realizing the Jordan canonical form.

Finally, let

$$
E^{1}=\left\{v \in E_{i \mu} \mid\left\langle L u_{p, q}^{(j)}, v\right\rangle=0, \forall 1 \leq j \leq j_{0}, 1 \leq p \leq l_{j}, 1 \leq q \leq k_{j}\right\}
$$

Much as in the invariance of $Z$ in the above, $J L\left(E^{1}\right) \subset E^{1}$. Since all the Jordan blocks are realized by

$$
\left\{u_{p, q}^{(j)}, 1 \leq j \leq j_{0}, 1 \leq p \leq l_{j}, 1 \leq q \leq k_{j}\right\}
$$

we have $E^{1} \subset \operatorname{ker}(J L-i \mu)$. This completes the proof.
Based on Proposition 2.2, we give the following result to be used later.
LEMMA 4.6. Let $J, L$ be real $n \times n$ matrices. Assume $J$ is anti-symmetric and $L$ is symmetric and nonsingular. Then there exists an invariant (under $J L$ ) subspace $W$ of $\mathbf{C}^{n}$ such that $\operatorname{dim} W=n^{-}(L)$ and $\left.\langle L \cdot, \cdot\rangle\right|_{W} \leq 0$.

Proof. For any purely imaginary eigenvalue $\lambda=i \mu \in i \mathbf{R}$, we start with the special basis of $E_{i \mu}$ given by Proposition 2.2 (as well as Remark 2.11). For each Jordan chain $\left\{u_{p, 1}^{(j)}, \cdots, u_{p, k_{j}}^{(j)}\right\}$ of even length, define the subspace $Z_{i \mu, j, p}=$ $\operatorname{span}\left\{u_{p, 1}^{(j)}, \cdots, u_{p, k_{j} / 2}^{(j)}\right\}$. For each Jordan chain $\left\{u_{p, 1}^{(j)}, \cdots, u_{p, k_{j}}^{(j)}\right\}$ of odd length $k_{j} \geq 1$, define the subspace

$$
Z_{i \mu, j, p}=\left\{\begin{array}{l}
\operatorname{span}\left\{u_{p, 1}^{(j)}, \cdots, u_{p,\left(k_{j}-1\right) / 2}^{(j)}\right\} \\
\operatorname{span}\left\{u_{p, 1}^{(j)}, \cdots, u_{p,\left(k_{j}+1\right) / 2}^{(j)}\right\}
\end{array}\right.
$$

$$
\begin{aligned}
& \text { if }\left\langle L u_{p,\left(k_{j}+1\right) / 2}^{(j)}, u_{p,\left(k_{j}+1\right) / 2}^{(j)}\right\rangle>0, \\
& \text { if }\left\langle L u_{p,\left(k_{j}+1\right) / 2}^{(j)}, u_{p,\left(k_{j}+1\right) / 2}^{(j)}\right\rangle<0 .
\end{aligned}
$$

Proposition 2.2 implies that $\langle L u, u\rangle \leq 0$ for all $u \in Z_{i \mu, j, p}$ defined above. For any eigenvalue $\lambda$ of $J L$ with $\operatorname{Re} \lambda>0$, recall $\langle L u, u\rangle=0$ for all $u \in E_{\lambda}$ by Lemma 3.3. Define

$$
Z_{i \mu}=\oplus_{j=0}^{k_{j}} \oplus_{p=1}^{l_{j}} Z_{i \mu, j, p}
$$

and

$$
W=\oplus_{\operatorname{Re} \lambda>0} E_{\lambda} \oplus_{i \mu \in \sigma(J L) \cap i \mathbf{R}} Z_{i \mu} .
$$

Then $\left.\langle L \cdot, \cdot\rangle\right|_{W} \leq 0$ since these subspaces are pairwise orthogonal in $\langle L \cdot, \cdot\rangle$. Moreover, $\operatorname{dim} W=n^{-}(L)$ due to the counting formula (2.13) and (2.16).

## CHAPTER 5

## Invariant subspaces

In this chapter, we study subspaces of $X$ invariant under $J L$, including both positive and negative results. As the first step to prove our main results, a nonpositive (with respect to $\langle L \cdot, \cdot\rangle$ ) invariant subspace of the maximal possible dimension $n^{-}(L)$ is derived in Section 5.1. The existence of such subspaces is not only useful for the linear dynamics, but also a rather interesting and delicate result as demonstrated in the discussions and examples in Section 5.2. Throughout this chapter, we work under the non-degeneracy assumption that (2.4) holds for $L$ which is equivalent to $L: X \rightarrow X^{*}$ is an isomorphism.

### 5.1. Maximal non-positive invariant subspaces

Theorem 5.1. In additional to hypotheses $(\boldsymbol{H}-3)$, assume $L$ satisfies the nondegeneracy assumption (2.4), then
(1) $\operatorname{dim} W \leq n^{-}(L)$ holds for any subspace $W \subset X$ satisfying $\langle L u, u\rangle \leq 0$ for any $u \in W$; and
(2) there exists a subspace $W \subset D(J L)$ such that

$$
\operatorname{dim} W=n^{-}(L), \quad J L(W) \subset W, \text { and }\langle L u, u\rangle \leq 0, \forall u \in W
$$

Remark 5.1. Though the theorem is stated for real Hilbert spaces, the same proof shows that it also holds for complex Hilbert space $X$ and Hermitian forms $L$ and $J$. Furthermore, the invariance of $W$ under $J L$ implies that $W \subset \cap_{k=1}^{\infty} D\left((J L)^{k}\right)$.

This theorem is basically equivalent to the classical Pontryagin invariant subspace theorem which is usually stated for a self-adjoint operator $A$ with resect to some indefinite quadratic form $\langle L \cdot, \cdot\rangle$ on $X$ with finitely many negative directions (i.e. $L$ satisfies (H2) with $n^{-}(L)<\infty$ and $\operatorname{ker} L=\{0\}$ ). It states that there exists a subspace $W \subset X$ such that $W$ is invariant under $A,\left.\langle L \cdot, \cdot\rangle\right|_{W} \leq 0$ and $\operatorname{dim} W=$ $n^{-}(L)$ (i.e. maximal non-positive dimension). Such theorems have been proved in the literature (e.g. see [28] [18] and the references therein). We believe that it will play a fundamental role in further studies of Hamiltonian systems and deserves more attention than it currently does. For the Hamiltonian PDE (2.1) considered in this paper, one important observation is that the operator $J L$ is anti-self-adjoint with respect to the inner product $\langle L \cdot, \cdot\rangle$. Since both anti-self-adjoint and self-adjoint operators are related to unitary operators by the Cayley transform, the Pontryagin invariant subspace theorems can be equivalently stated for unitary, self-adjoint or anti-self-adjoint cases. By Lemma 3.1, $e^{t J L}$ is unitary in $\langle L \cdot, \cdot\rangle$. But to study the eigenvalues of $J L$ more directly, we still use Cayley transform to relate $J L$ to an unitary operator and then apply the Pontryagin invariant subspace theorem. For the sake of completeness, in the following we outline a proof of Theorem 5.1 by the arguments given in [18] for the proof of Pontryagin invariant subspace theorem via
unitary operators which is based on compactness and fixed point theorems (see also [25] [48]).

We also give another more constructive proof of Theorem 5.1, by using the Hamiltonian structure of (2.1) and Galerkin approximation. It provides more information about the invariant subspace $W$.

Proof. The assumption (2.4) is equivalent to $\operatorname{ker} L=\{0\}$. The first statement of the Theorem follows by the same proof of Lemma 12.1. Below we give two different proofs of the construction of the invariant subspace $W$ in the second part of the Theorem.

Proof (\#1.) Here we sketch a proof of Theorem 5.1 by using the arguments in [18]. Let $X_{ \pm} \subset X$ be given by Lemma 12.4. Assumptions (H2-3) and (2.4) ensure that

$$
X=X_{-} \oplus X_{+}, \quad X^{*}=\tilde{X}_{-}^{*} \oplus \tilde{X}_{+}^{*}, \quad \pm\langle L u, u\rangle \geq \delta\|u\|^{2}, \quad \forall u \in X_{ \pm}
$$

where $\tilde{X}_{ \pm}^{*}=P_{ \pm}^{*} X_{ \pm}^{*}$ and $P_{ \pm}$are the associated projections. As in the proof of Lemma 12.5, let $i_{X_{ \pm}}: X_{ \pm} \rightarrow X$ be the embedding and

$$
L_{ \pm}= \pm P_{ \pm}^{*} i_{X_{ \pm}}^{*} L i_{X_{ \pm}} P_{ \pm}, \quad(u, v)_{L} \triangleq\left\langle\left(L_{+}+L_{-}\right) u, v\right\rangle
$$

There exists $\delta>0$ such that $\left\langle L_{ \pm} u, u\right\rangle \geq \delta\|u\|^{2}$, for all $u \in X_{ \pm}$and the quadratic form $(\cdot, \cdot)_{L}$ induces an equivalent norm $|u|_{L} \triangleq \sqrt{(u, v)_{L}}$ on $X$. We denote the Hilbert space $\left(X,\left\langle\left(L_{+}+L_{-}\right) \cdot, \cdot\right\rangle\right)$ by $X_{L}$.

Step 1. It is clear that $J\left(L_{+}+L_{-}\right)$is an anti-self-adjoint operator on $X_{L}$ and $J L_{-}$is a bounded linear operator of finite rank on $X_{L}$ as $L_{-} X_{-}=P_{-}^{*} X_{-}^{*} \subset D(J)$. Writing $J L=J\left(L_{+}+L_{-}\right)-2 J L_{-}$, we obtain that there exists $a>0$ such that $\alpha \notin \sigma(J L)$ if $|\operatorname{Re} \alpha| \geq a$. Let

$$
T=(J L+a)(J L-a)^{-1}, \text { then }\langle L T u, T v\rangle=\langle L u, v\rangle, \forall u, v \in X
$$

through straightforward calculation using $J=-J^{*}$ and that $L$ is bounded and symmetric. In some sense, $J L$ is anti-self-adjoint with respect to the quadratic form $\langle L \cdot, \cdot\rangle$ and thus $T$ is formally the Cayley transformation.

Step 2. Let $X_{L \pm}$ be the subspaces $X_{ \pm}$equipped with the inner product $(\cdot, \cdot)_{L}$ which is equivalent to $\pm\left\langle L_{X_{ \pm}} \cdot, \cdot\right\rangle$ on $X_{ \pm}$, where $L_{X_{ \pm}}$is defined in (12.1). One may prove (see Lemma 3.6 in [18]) that a subspace $W \subset X_{L}$ satisfies $\operatorname{dim} W=$ $n^{-}(L)$ and $\langle L u, u\rangle \leq 0$ for all $u \in W$ if and only if $W$ is the graph of a bounded linear operator $S: X_{L-} \rightarrow X_{L_{+}}$with operator norm $|S| \leq 1$. Denote this set of operators, i.e. the unit ball of $L\left(X_{L-}, X_{L+}\right)$, by $B_{1}\left(X_{L-}, X_{L+}\right)$. This proves the first statement.

Step 3. For any $S \in B_{1}\left(X_{L-}, X_{L+}\right)$, since $T$ preserves the quadratic form $\langle L \cdot, \cdot\rangle$, one may show that $T(\operatorname{graph}(S))$ is still the graph of some $S^{\prime} \in B_{1}\left(X_{L-}, X_{L+}\right)$. Hence we define a transformation $\mathcal{T}$ on $B_{1}\left(X_{L-}, X_{L+}\right)$ as

$$
\operatorname{graph}(\mathcal{T}(S))=T(\operatorname{graph}(S))
$$

Step 4. The space of bounded operators $L\left(X_{L-}, X_{L+}\right)$ equipped with the weak topology is a locally convex topological vector space. Since $X_{L-}$ is finite dimensional, the unit ball $B_{1}\left(X_{L-}, X_{L+}\right)$ is convex and compact under the weak topology. Using the boundedness and the finite dimensionality of $X_{L-}$, one may prove (see [18] for details) that $\mathcal{T}$ is continuous under the weak topology. According to the Tychonoff fixed point theorem (sometimes referred as the SchauderTychonoff fixed point theorem, see [76]), $\mathcal{T}$ has a fixed point $S \in B_{1}\left(X_{L-}, X_{L+}\right)$.

Let $W=\operatorname{graph}(S)$ and thus $T(W) \subset W$. According to the definition of $T$, we have

$$
(J L-a)^{-1}=\frac{1}{2 a}(T-I)
$$

which implies that $W$ is invariant under $(J L-a)^{-1}$. As $W$ is finite dimensional and $(J L-a)^{-1}$ is bounded and injective, it is clear that $W=(J L-a)^{-1} W \subset D(J L)$ and thus $J L(W) \subset W$.

Alternative proof (\#2) of Theorem 5.1 via Galerkin approximation on separable $X$. On the one hand, the above proof given in [18] is elegant and is based on fixed point theorems involving compactness, which does not yield much detailed information of the invariant subspace W . On the other hand, clearly Theorem 5.1 is a generalization into Hilbert spaces of Lemma 4.6 whose constructive proof provides more explicit information of the invariant subspaces. In fact, assuming $X$ is separable, in the rest of this chapter we give an alternative proof of Theorem 5.1 based on Lemma 4.6.

Denote

$$
\begin{equation*}
[\cdot, \cdot]=\langle L \cdot, \cdot\rangle \text { on } X \tag{5.1}
\end{equation*}
$$

Let $X_{ \pm}$be the same subspaces of $X$ chosen as in the above proof $\# 1$ (as well as in the proof of Proposition 12.1). We will study the eigenvalues of $J L$ by a Galerkin approximation. Choose an orthogonal (with respect to $[\cdot, \cdot]$ ) basis $\left\{\xi_{k}\right\}_{k=1}^{\infty}$ of $X$ such that $\xi_{k} \in D(J L)$,

$$
X_{-}=\operatorname{span}\left\{\xi_{1}, \cdots, \xi_{n^{-}(L)}\right\}, X_{+}=\overline{\operatorname{span}\left\{\xi_{k}\right\}_{k=n^{-}(L)+1}^{\infty}},
$$

and $\left[\xi_{k}, \xi_{j}\right]=0$ if $k \neq j ;\left[\xi_{j}, \xi_{j}\right]=-1$ if $1 \leq j \leq n^{-}(L) ;\left[\xi_{j}, \xi_{j}\right]=1$ if $j \geq$ $n^{-}(L)+1$. For each $n>n^{-}(L)$, define $X^{(n)}=\operatorname{span}\left\{\xi_{1}, \cdots, \xi_{n}\right\}$ and denote $\pi^{n}$ be the orthogonal projection with respect to the quadratic form $[\cdot, \cdot]$ from $X$ to $X^{(n)}$.

Let $X, J L$, and $[\cdot, \cdot]$ (as a Hermitian symmetric form) also denote their complexifications as in Chapter 3. Still $\left\{\xi_{1}, \xi_{2}, \ldots\right\}$ form a basis of the complexified $X$. Define the operator $F^{(n)}: X^{(n)} \rightarrow X^{(n)}$ by

$$
F^{(n)} v=\pi^{n} J L v
$$

Notice that, for $j, k \leq n$,

$$
\left[F^{(n)} \xi_{k}, \xi_{j}\right]=\left[\pi^{n} J L \xi_{k}, \xi_{j}\right]=\left\langle L J L \xi_{k}, \xi_{j}\right\rangle=\left\langle L \xi_{j}, J L \xi_{k}\right\rangle \triangleq\left(J^{(n)}\right)_{j k}
$$

where the $n \times n$ matrix $\left(J^{(n)}\right)$ is real and anti-symmetric. Let $v=\sum_{j=1}^{n} y_{j} \xi_{j} \in X^{(n)}$ and denote $\vec{y}^{(n)}=\left(y_{1}, \cdots, y_{n}\right)^{T}$ and the $n \times n$ matrix

$$
H^{(n)}=\left(\left[\xi_{k}, \xi_{j}\right]\right)=\operatorname{diag}[\underbrace{-1, \cdots,-1}_{1 \text { to } n^{-(L)}}, \underbrace{1, \cdots, 1}_{n^{-}(L)+1 \text { to } n}] .
$$

Then $F^{(n)} v=\sum_{k=1}^{n} a_{k} \xi_{k}$, where

$$
\vec{a}^{(n)}=\left(a_{1}, \cdots, a_{n}\right)^{T}=H^{(n)} J^{(n)} \vec{y}^{(n)}
$$

So the eigenvalue problem $F^{(n)}(v)=\lambda v$ is equivalent to

$$
\begin{equation*}
H^{(n)} J^{(n)} \vec{y}^{(n)}=\lambda \vec{y}^{(n)}, \tag{5.2}
\end{equation*}
$$

Let $\vec{z}^{(n)}=H^{(n)} \vec{y}^{(n)}$, then the eigenvalue problem (5.2) becomes

$$
J^{(n)} H^{(n)} \vec{z}^{(n)}=\lambda \vec{z}^{(n)}
$$

For any $n \geq n^{-}(L)$, since $n^{-}\left(H^{(n)}\right)=n^{-}(L)$, by Lemma 4.6, there exists a subspace $Z^{(n)} \subset \mathbf{C}^{n}$ of dimension $n^{-}(L)$, such that $Z^{(n)}$ is invariant under $J^{(n)} H^{(n)}$ and $\left\langle H^{(n)} z, z\right\rangle \leq 0$ for any $z \in Z^{(n)}$. Define $Y^{(n)}=H^{(n)} Z^{(n)}$ and

$$
W^{(n)}=\left\{\sum_{j=1}^{n} y_{j} \xi_{j} \mid\left(y_{1}, \cdots, y_{n}\right)^{T} \in Y^{(n)}\right\}
$$

Then $W^{(n)}$ is invariant under the linear mapping $F^{(n)}$, $\operatorname{dim}\left(W^{(n)}\right)=n^{-}(L)$ and the quadratic functions

$$
\left.\left.\langle L \cdot, \cdot\rangle\right|_{W^{(n)}}=\left\langle H^{(n)}, \cdot\right\rangle\right\rangle\left.\right|_{Y^{(n)}}=\left.\left\langle H^{(n)} \cdot, \cdot\right\rangle\right|_{Z^{(n)}} \leq 0
$$

As in the proof $(\# 1)$ above, denote $P_{ \pm}: X \rightarrow X_{ \pm}$to be the projection operators with $\operatorname{ker} P_{ \pm}=X_{\mp}$. Since the definitions of $X_{+}$and $W^{(n)}$ imply $W^{(n)} \cap X_{+}=\{0\}$, it holds that $P_{-}\left(W^{(n)}\right)=X_{-}$. So we can choose a basis $\left\{w_{1}^{(n)}, \cdots, w_{n^{-}(L)}^{(n)}\right\}$ of $W^{(n)}$ such that $w_{j}^{(n)}=\xi_{j}+w_{j+}^{(n)}$ with $w_{j+}^{(n)} \in X_{+}$. For each $j \leq n^{-}(L)$, since

$$
0 \geq\left\langle L w_{j}^{(n)}, w_{j}^{(n)}\right\rangle=\left\langle L \xi_{j}^{(n)}, \xi_{j}^{(n)}\right\rangle+\left\langle L w_{j+}^{(n)}, w_{j+}^{(n)}\right\rangle \geq-1+\delta_{0}\left\|w_{j+}^{(n)}\right\|^{2}
$$

so $\left\|w_{j}^{(n)}\right\| \leq C$ for some constant $C$ independent of $j$ and $n$. Therefore, as $n \rightarrow \infty$, subject to a subsequence, we have $w_{j}^{(n)} \rightharpoonup w_{j}^{\infty} \in X$ weakly and $P_{-}\left(w_{j}^{\infty}\right)=\xi_{j}$. The subspace $W^{\infty}=\operatorname{span}\left\{w_{j}^{\infty}\right\}_{j=1}^{n^{-}(L)}$ is of dimension $n^{-}(L)$ since $P_{-}\left(W^{\infty}\right)=X_{-}$.

We now show that: i) $W^{\infty}$ is invariant under the operator $J L$ and ii) $\langle L u, u\rangle \leq 0$ for any $u \in W^{\infty}$. To prove i), first note that since $W^{(n)}$ is invariant under $F^{(n)}$, we have

$$
F^{(n)} w_{k}^{(n)}=\sum_{j=1}^{n^{-}(L)} a_{k j}^{(n)} w_{j}^{(n)}, \quad a_{i j}^{(n)} \in \mathbf{C} .
$$

For any integer $l \in \mathbf{N}$ and a fixed $w \in X^{(l)}$, when $n \geq l$,

$$
\begin{align*}
\sum_{j=1}^{n^{-}(L)} a_{k j}^{(n)}\left[w_{j}^{(n)}, w\right] & =\left[F^{(n)} w_{k}^{(n)}, w\right]=\left\langle L J L w_{k}^{(n)}, w\right\rangle  \tag{5.3}\\
& =-\left\langle L w_{k}^{(n)}, J L w\right\rangle=-\left[w_{k}^{(n)}, J L w\right]
\end{align*}
$$

We claim that $\left\{a_{i j}^{(n)}\right\}$ is uniformly bounded for $1 \leq k, j \leq n^{-}(L)$ and $n>n^{-}(L)$. Suppose otherwise, there exists $1 \leq k_{0}, j_{0} \leq n^{-}(L)$ and a subsequence $\left\{n_{m}\right\} \rightarrow \infty$, such that, for all $j \leq n^{-}(L)$,

$$
\left|a_{k_{0} j_{0}}^{\left(n_{m}\right)}\right|=\max _{1 \leq j \leq n^{-}(L)}\left\{\left|a_{k_{0} j}^{\left(n_{m}\right)}\right|\right\} \rightarrow \infty \text { and } \forall j, c_{k_{0}, j}=\lim _{m \rightarrow \infty} a_{k_{0} j}^{\left(n_{m}\right)} / a_{k_{0} j_{0}}^{\left(n_{m}\right)} \text { exists }
$$

Then from (5.3), we get

$$
\sum_{j=1}^{n^{-}(L)} \frac{a_{k_{0} j}^{\left(n_{m}\right)}}{a_{k_{0} j_{0}}^{\left(n_{m}\right)}}\left[w_{j}^{\left(n_{m}\right)}, w\right]=-\frac{1}{a_{k_{0} j_{0}}^{\left(n_{m}\right)}}\left[w_{k_{0}}^{\left(n_{m}\right)}, J L w\right]
$$

and letting $m \rightarrow \infty$, we obtain

$$
\begin{equation*}
\sum_{j=1}^{n^{-}(L)} c_{k_{0}, j}\left[w_{j}^{\infty}, w\right]=0 \tag{5.4}
\end{equation*}
$$

where in particular we also notice $\left|c_{k_{0}, j_{0}}\right|=1$. By a density argument, the identity (5.4) holds also for any $w \in X$. Therefore, $\sum_{j=1}^{n^{-}(L)} c_{k_{0}, j} w_{j}^{\infty}=0$ by the nondegeneracy of $[\cdot, \cdot]$. This is in contradiction to the independency of $\left\{w_{k}^{\infty}\right\}$. So $\left\{a_{k j}^{(n)}\right\}$ is uniformly bounded. Let $n \rightarrow \infty$ in (5.3), subject to a subsequence, we obtain

$$
\sum_{j=1}^{n^{-}(L)} a_{k j}^{\infty}\left[w_{j}^{\infty}, w\right]=-\left[w_{j}^{\infty}, J L w\right]=\left[J L w_{k}^{\infty}, w\right], \text { where } a_{k j}^{\infty}=\lim _{n \rightarrow \infty} a_{k j}^{(n)}
$$

By a density argument again the above equality is also true for any $w \in X$, which implies

$$
J L\left(w_{k}^{\infty}\right)=\sum_{j=1}^{n^{-}(L)} a_{k j}^{\infty} w_{j}^{\infty}
$$

So $W^{\infty}$ is invariant under $J L$.
Now we prove the above claim ii), that is, $\left.\langle L \cdot, \cdot\rangle\right|_{W^{\infty}} \leq 0$. For any

$$
u=\sum_{j=1}^{n^{-}(L)} c_{j} w_{j}^{\infty} \in W^{\infty}
$$

denote

$$
u^{(n)}=\sum_{j=1}^{n^{-}(L)} c_{j} w_{j}^{(n)} \in W^{(n)}
$$

Clearly, $u^{(n)} \rightharpoonup u$ weakly in $X$ and $\left\langle L u^{(n)}, u^{(n)}\right\rangle \leq 0$, which converges subject to a subsequence. Since

$$
\lim _{n \rightarrow \infty}\left\langle L P_{-} u^{(n)}, P_{-} u^{(n)}\right\rangle=\left\langle L P_{-} u, P_{-} u\right\rangle
$$

which is due to $P_{-} w_{j}^{\infty}=\xi_{j}$ and therefore $P_{-} u^{(n)} \rightarrow P_{-} u$ strongly in $X$, and

$$
\lim _{n \rightarrow \infty}\left\langle L P_{+} u^{(n)}, P_{+} u^{(n)}\right\rangle \geq\left\langle L P_{+} u, P_{+} u\right\rangle
$$

as $\langle L x, x\rangle^{\frac{1}{2}}$ is a norm on $X_{+}$. Therefore,

$$
\begin{align*}
0 & \geq \lim _{n \rightarrow \infty}\left\langle L u^{(n)}, u^{(n)}\right\rangle=\left\langle L P_{-} u, P_{-} u\right\rangle+\lim _{n \rightarrow \infty}\left\langle L P_{+} u^{(n)}, P_{+} u^{(n)}\right\rangle  \tag{5.5}\\
& \geq\left\langle L P_{-} u, P_{-} u\right\rangle+\left\langle L P_{+} u, P_{+} u\right\rangle=\langle L u, u\rangle
\end{align*}
$$

This complete the proof of claim ii) and thus the proof of Theorem 5.1 under the separable assumption on $X$.

### 5.2. Further discussions on invariant subspaces and invariant decompositions

Continuous dependence of invariant subspaces on $J L$. In perturbation problems, the operator $J L$ may depend on a perturbation parameter $\epsilon$. One would naturally wish that a family $W_{\epsilon}$ of non-positive invariant subspaces of dimension $n^{-}(L)$ may be found depending on $\epsilon$ at least continuously. However, this turns out to be impossible in general, even if $L$ is assumed to be non-degenerate. See an example in Section 8.3.

Invariant splitting, I. In the presence of $W$ invariant under $J L$ with $\operatorname{dim} W=$ $n^{-}(L)$, it is natural to ask whether it is possible to make it into an invariant (under $J L)$ decomposition of $X$, i.e. whether there exist such $W$ and a codim- $n^{-}(L)$ invariant subspace $W_{1} \subset X$ such that $X=W \oplus W_{1}$. This is usually not possible as in the following example

$$
J=\left(\begin{array}{cccc}
0 & -1 & 1 & 0 \\
1 & 0 & 0 & 1 \\
-1 & 0 & 0 & 0 \\
0 & -1 & 0 & 0
\end{array}\right), L=\left(\begin{array}{cccc}
0 & 0 & 0 & 1 \\
0 & 0 & -1 & 0 \\
0 & -1 & 0 & 0 \\
1 & 0 & 0 & 0
\end{array}\right), J L=\left(\begin{array}{cccc}
0 & -1 & 1 & 0 \\
1 & 0 & 0 & 1 \\
0 & 0 & 0 & -1 \\
0 & 0 & 1 & 0
\end{array}\right) .
$$

Here $n^{-}(L)=2$ and the only eigenvalues are $\sigma(J L)=\{ \pm i\}$. The only possible non-positive 2 -dim invariant subspace, where the eigenvalues of the restriction of $J L$ are contained in $\sigma(J L)$, has to be the geometric kernel of $\pm i$ and thus $W=\left\{x_{3}=x_{4}=0\right\}$. There does not exist any 2-dim invariant subspace $W_{1}$ such that $\mathbf{R}^{4}=W \oplus W_{1}$ since the restriction of $J L$ on $W_{1}$ has to have eigenvectors of $\pm i$ as well.

Invariant Splitting, II. In light of Lemma 3.2,

$$
W^{\perp_{L}}=\{u \in X \mid\langle L u, v\rangle=0, \forall v \in W\}
$$

is invariant under $e^{t J L}$. While one may wish $X=W \oplus W^{\perp_{L}}$, the only obstacle is that $W^{\perp_{L}}$ may intersect $W$ nontrivially as $L_{W}$, as defined in (12.1), may be degenerate as in the above example. A more natural question is whether it is possible to enlarge $W$ to some closed $\tilde{W} \supset W$ such that

$$
\operatorname{dim} \tilde{W}<\infty, \quad J L(\tilde{W}) \subset \tilde{W}, \quad \text { and } L_{\tilde{W}} \text { an isomorphism. }
$$

If so, Lemmas 12.2 and 3.2 would imply

$$
\tilde{W}^{\perp_{L}}=\{u \in X \mid\langle L u, v\rangle=0, \forall v \in \tilde{W}\}
$$

is invariant under $e^{t J L}$ and $X=\tilde{W} \oplus \tilde{W}^{\perp_{L}}$. Moreover, $L_{\tilde{W}^{\perp_{L}}}$ is positive definite due to Theorem 5.1 and thus $e^{t J L}$ is stable on $\tilde{W}^{\perp_{L}}$. Consequently all the index counting and stability analysis related to $e^{t J L}$ can be reduced to the finite dimensional $\tilde{W}$, which has been analyzed in Chapter 4. For example, we would have a counting theorem like Theorem 2.3, in particular with $k_{0}^{\leq 0}$ and $k_{i}^{\leq 0}$ replaced by $k_{0}^{-}$and $k_{i}^{-}$, respectively.

Unfortunately the above splitting is not always possible either, as can be seen from a counterexample in Section 8.4. In Proposition 2.4, we give conditions to get such a decomposition.

## CHAPTER 6

## Structural decomposition

In this chapter, we prove Theorem 2.1 and Corollary 2.1 on the decomposition of $X$. Our first step to decompose $X$ is the following proposition based on the invariant subspace Theorem 5.1.

Proposition 6.1. In addition to (H1-H3), assume ker $L=\{0\}$. There exist closed subspaces $Y_{j}, j=1,2,3,4$, such that $X=\oplus_{j=1}^{4} Y_{j}$ and

$$
\begin{equation*}
\operatorname{dim} Y_{1}=\operatorname{dim} Y_{4}=n^{-}(L)-\operatorname{dim} Y_{2}<\infty, \quad Y_{1,2,4} \subset \cap_{k=1}^{\infty} D\left((J L)^{k}\right) \tag{6.1}
\end{equation*}
$$

and accordingly the linear operator $J L$ and the quadratic form $\langle L \cdot, \cdot\rangle$ take the block forms

$$
J L \longleftrightarrow\left(\begin{array}{cccc}
\tilde{A}_{1} & \tilde{A}_{12} & \tilde{A}_{13} & \tilde{A}_{14} \\
0 & \tilde{A}_{2} & 0 & \tilde{A}_{24} \\
0 & 0 & \tilde{A}_{3} & \tilde{A}_{34} \\
0 & 0 & 0 & \tilde{A}_{4}
\end{array}\right), \quad L \longleftrightarrow\left(\begin{array}{cccc}
0 & 0 & 0 & \tilde{B} \\
0 & L_{Y_{2}} & 0 & 0 \\
0 & 0 & L_{Y_{3}} & 0 \\
\tilde{B}^{*} & 0 & 0 & 0
\end{array}\right)
$$

Here $\tilde{B}: Y_{4} \rightarrow Y_{1}^{*}$ is an isomorphism and the quadratic forms $L_{Y_{2}} \leq-\delta_{0}$ and $L_{Y_{3}} \geq \delta_{0}$ for some $\delta_{0}>0$. Moreover, $\tilde{A}_{3}: D\left(\tilde{A}_{3}\right)=Y_{3} \cap D(J L) \rightarrow Y_{3}$ is closed, while all other blocks are bounded operators. The operators $\tilde{A}_{2,3}$ are anti-self-adjoint with respect to the equivalent inner product $\mp\left\langle L_{Y_{2,3}}, \cdot\right\rangle$.

Before we give the proof the proposition, we would like to make two remarks. Firstly we observe that $(J L)^{k}$ takes the same blockwise form as the above one of $J L$. Secondly, the bounded operator $\tilde{A}_{13}$ should be understood as the closure of $\left.P_{1} J L\right|_{Y_{3}}$, which may not be closed or everywhere defined itself. Here $P_{j}: X \rightarrow Y_{j}$ is the projection to $Y_{j}, j=1,2,3,4$, according to the decomposition.

Proof. Theorem 5.1 states that there exists $W \subset D\left((J L)^{k}\right)$ such that $\operatorname{dim} W=$ $n^{-}(L), J L(W) \subset W$, and $\langle L u, u\rangle \leq 0$ for all $u \in W$. Let $Y_{1}=W \cap W^{\perp_{L}}\left(\perp_{L}\right.$ defined as in Lemma 12.2), $\tilde{Y}_{2} \subset W$, and $\tilde{Y}_{3} \subset W^{\perp_{L}}$ be closed subspaces such that $W=Y_{1} \oplus \tilde{Y}_{2}$ and $W^{\perp_{L}}=Y_{1} \oplus \tilde{Y}_{3}$. Recall the notation $L_{Y}$ as defined in (12.1) for any closed subspace $Y$.

Claim. $L_{Y_{1}}=0$ and there exists $\delta_{0}>0$ such that $L_{\tilde{Y}_{2}} \leq-\delta_{0}$ and $L_{\tilde{Y}_{3}} \geq \delta_{0}$.
In fact, since the quadratic form $\langle L u, u\rangle \leq 0$, for all $u \in W$, the variational principle yields that $u \in W$ satisfies $\langle L u, u\rangle=0$ if and only if $\langle L u, v\rangle=0$ for all $v \in W$, or equivalently $u \in W \cap W^{\perp_{L}}=Y_{1}$. Therefore, $\langle L u, u\rangle<0$ for any $u \in \tilde{Y}_{2} \backslash\{0\}$ which along with $\operatorname{dim} \tilde{Y}_{2}<\infty$ implies $L_{\tilde{Y}_{2}} \leq-\delta_{0}$ for some $\delta_{0}>0$.

If there exists $u \in W^{\perp_{L}} \backslash W$ satisfying $\langle L u, u\rangle \leq 0$, the definition of $W^{\perp_{L}}$ would imply $\langle L u, v\rangle \leq 0$ for all $v \in \tilde{W}=W \oplus \mathbf{R} u$ and $\operatorname{dim} \tilde{W}=n^{-}(L)+1$. This would contradict Theorem 5.1 and thus we obtain $\langle L u, u\rangle>0$ for all $u \in \tilde{Y}_{3} \backslash\{0\}$.

Consequently, Lemma 12.2 implies that $L_{\tilde{Y}_{3}} \geq \delta_{0}$ for some $\delta_{0}>0$ and the claim is proved.

Since $L$ is assumed to non-degenerate, it is easy to see codim- $\left(W+W^{\perp_{L}}\right)=$ $\operatorname{dim} Y_{1}<\infty$. Let $\tilde{Y}_{4}$ be a subspace such that $X=\left(W+W^{\perp_{L}}\right) \oplus \tilde{Y}_{4}=Y_{1} \oplus \tilde{Y}_{2} \oplus \tilde{Y}_{3} \oplus \tilde{Y}_{4}$ and $\tilde{Y}_{4} \subset \cap_{k=1}^{\infty} D\left((J L)^{k}\right)$, which is possible as $\cap_{k=1}^{\infty} D\left((J L)^{k}\right)$ is dense and $\operatorname{dim} \tilde{Y}_{4}=$ $\operatorname{dim} Y_{1}<\infty$. With respect to this decomposition, $L$ takes the form

$$
L \longleftrightarrow\left(\begin{array}{cccc}
0 & 0 & 0 & \tilde{B}_{41}^{*} \\
0 & L_{\tilde{Y}_{2}} & 0 & \tilde{B}_{42}^{*} \\
0 & 0 & L_{\tilde{Y}_{3}} & \tilde{B}_{43}^{*} \\
\tilde{B}_{41} & \tilde{B}_{42} & \tilde{B}_{43} & L_{\tilde{Y}_{4}}
\end{array}\right)
$$

The non-degeneracy of $L$ implies that $\tilde{B}_{41}=i_{\tilde{Y}_{4}}^{*} L i_{Y_{1}}: Y_{1} \rightarrow \tilde{Y}_{4}^{*}$ is an isomorphism. Let

$$
S_{4}=-\frac{1}{2} \tilde{B}_{41}^{-1} L_{\tilde{Y}_{4}}: \tilde{Y}_{4} \rightarrow Y_{1}, \quad S_{j}=-\tilde{B}_{41}^{-1} \tilde{B}_{4 j}: \tilde{Y}_{j} \rightarrow Y_{1}, j=2,3
$$

For any $u, v \in \tilde{Y}_{4}$, we have

$$
\begin{aligned}
\left\langle L\left(u+S_{4} u\right), v+S_{4} v\right\rangle & =\left\langle L_{\tilde{Y}_{4}} u, v\right\rangle+\left\langle\left(L S_{4}+\left(L S_{4}\right)^{*}\right) u, v\right\rangle \\
& =\left\langle L_{\tilde{Y}_{4}} u, v\right\rangle+\left\langle\left(\tilde{B}_{41} S_{4}+\left(\tilde{B}_{41} S_{4}\right)^{*}\right) u, v\right\rangle=0 .
\end{aligned}
$$

Similarly, for any $u \in \tilde{Y}_{j}, j=2,3$, and $v \in \tilde{Y}_{4}$,

$$
\left\langle L\left(u+S_{j} u\right), v+S_{4} v\right\rangle=\left\langle\tilde{B}_{4 j} u, v\right\rangle+\left\langle\tilde{B}_{41} S_{j} u, v\right\rangle=0
$$

Let $Y_{j}=\left(I+S_{j}\right) \tilde{Y}_{j}$. Clearly, it still holds $X=\oplus_{j=1}^{4} Y_{4}$. Moreover, $Y_{1,2,4} \subset$ $\cap_{k=1}^{\infty} D\left((J L)^{k}\right)$, the dimension relationship in (6.1) holds, and in this decomposition $L$ takes the desired form as in the statement of the proposition. Due to $W=Y_{1} \oplus Y_{2}$ and $W^{\perp_{L}}=Y_{1} \oplus Y_{3}$, the same claim as above implies the uniform positivity of $-L_{Y_{2}}$ and $L_{Y_{3}}$. The non-degeneracy of $\tilde{B}$ follows from the non-degeneracy assumption of $L$.

The invariance of $W$, and thus the invariance of $W^{\perp_{L}}$ due to Lemma 3.2, yields the desired form of $J L$. The properties that $\tilde{A}_{2,3}$ are anti-self-adjoint with respect to $\left\langle L_{Y_{2,3}} \cdot, \cdot\right\rangle$ and the boundedness of other blocks can be proved by applying Lemma 12.3 repeatedly to the splitting based on $X=\left(Y_{1} \oplus Y_{4}\right) \oplus\left(Y_{2} \oplus Y_{3}\right)$.

The following general functional analysis lemma on invariant subspaces will be used several times in the rest of the paper.

Lemma 6.1. Let $Z$ be a Banach space and $Z_{1,2} \subset Z$ be closed subspaces such that $Z=Z_{1} \oplus Z_{2}$. Suppose $A$ is a linear operator on $X$ which, in the above splitting, takes the form $\left(\begin{array}{cc}A_{1} & A_{12} \\ 0 & A_{2}\end{array}\right)$, such that

- $A_{1,2}: Z_{1,2} \supset D\left(A_{1,2}\right) \rightarrow Z_{1,2}$ are densely defined closed operators, one of which and $A_{12}: Z_{2} \rightarrow Z_{1}$ is bounded and
- $\sigma\left(A_{1}\right) \cap \sigma\left(A_{2}\right)=\emptyset$,
then there exists a bounded operator $S: Z_{2} \rightarrow Z_{1}$ such that
(1) $S Z_{2} \subset D\left(A_{1}\right)$ and
(2) $A\left(\tilde{Z}_{2} \cap D(A)\right) \subset \tilde{Z}_{2}$, where $\tilde{Z}_{2}=(I+S) Z_{2}=\left\{z_{2}+S\left(z_{2}\right) \mid z_{2} \in Z_{2}\right\}$.

REMARK 6.1. Clearly, the above properties also imply $D(A) \cap \tilde{Z}_{2}=(I+S) D\left(A_{2}\right)$ is dense in the closed subspace $\tilde{Z}_{2}$ and $\left.A\right|_{\tilde{Z}_{2}}: D(A) \cap \tilde{Z}_{2} \rightarrow \tilde{Z}_{2}$ is a closed operator. By using the splitting $Z=Z_{1} \oplus \tilde{Z}_{2}$, $A$ is block diagonalized into $\operatorname{diag}\left(A_{1}, A_{2}\right)$. Moreover, if $A_{2}$ is bounded, then the closed graph theorem implies that $\left.A\right|_{\tilde{Z}_{2}}$ is also bounded.

The proof of this lemma may be found in some standard functional analysis textbook. For the sake of completeness we also give a proof here.

Proof. Let us first consider the case when $A_{2}$ is bounded. Since $\sigma\left(A_{2}\right)$ is compact and $\sigma\left(A_{2}\right) \cap \sigma\left(A_{1}\right)=\emptyset$, there exists an open subset $\Omega \subset \mathbf{C}$ with compact closure and smooth boundary $\Gamma=\partial \Omega$ such that $\sigma\left(A_{2}\right) \subset \Omega \subset \bar{\Omega} \subset \mathbf{C} \backslash \sigma\left(A_{1}\right)$. We have

$$
\frac{1}{2 \pi i} \oint_{\Gamma}\left(\lambda-A_{1}\right)^{-1} d \lambda=0, \quad \frac{1}{2 \pi i} \oint_{\Gamma}\left(\lambda-A_{2}\right)^{-1} d \lambda=I
$$

Define

$$
S=\frac{1}{2 \pi i} \oint_{\Gamma} T(\lambda) d \lambda, \text { where } T(\lambda)=\left(A_{1}-\lambda\right)^{-1} A_{12}\left(A_{2}-\lambda\right)^{-1}
$$

Since $\left(A_{j}-\lambda\right)^{-1}, j=1,2$, is analytic from $\mathbf{C} \backslash \sigma\left(A_{j}\right)$ to $L\left(Z_{j}\right)$, it is clear that $S: Z_{2} \rightarrow Z_{1}$ is bounded. In particular, observing $T(\lambda) z \in D\left(A_{1}\right)$ for any $z \in Z_{2}$, one may verify

$$
\begin{equation*}
T(\lambda) A_{2} z-A_{1} T(\lambda) z=\left(A_{1}-\lambda\right)^{-1} A_{12} z-A_{12}\left(A_{2}-\lambda\right)^{-1} z \triangleq \tilde{T}(\lambda) z \tag{6.2}
\end{equation*}
$$

where $\tilde{T}(\lambda) \in L\left(Z_{2}, Z_{1}\right)$ is also analytic in $\lambda$.
We first show that $S z \in D\left(A_{1}\right)$ for any $z \in Z_{2}$. In fact, let $S_{n}, n \in \mathbf{N}$, be the values of a sequence of Riemann sums of the integral defining $S$, such that $S_{n} \rightarrow S$. Clearly, the discrete Riemann sums satisfy $S_{n} z \in D\left(A_{1}\right)$ and along with (6.2) we obtain that

$$
S_{n} A_{2} z-A_{1} S_{n} z=\tilde{T}_{n} z \rightarrow \frac{1}{2 \pi i} \oint_{\Gamma} \tilde{T}(\lambda) z d \lambda=A_{12} z
$$

where $\tilde{T}_{n} z$ is the corresponding Riemann sum of the integral on the right side. Therefore, we obtain from the closedness of $A_{1}$ that $S z \in D\left(A_{1}\right)$ and

$$
\begin{equation*}
S A_{2}-A_{1} S=A_{12} \tag{6.3}
\end{equation*}
$$

From this equation it is straightforward to verify, for any $z \in Z_{2}$,

$$
\begin{equation*}
A(z+S z)=A_{2} z+S A_{2} z \tag{6.4}
\end{equation*}
$$

In the other case where $A_{1}$ is bounded, the proof is similar. In fact, let $\Omega \subset \mathbf{C}$ be an open subset with compact closure and smooth boundary $\Gamma=\partial \Omega$ such that $\sigma\left(A_{1}\right) \subset \Omega \subset \bar{\Omega} \subset \mathbf{C} \backslash \sigma\left(A_{2}\right)$. Define

$$
S=-\frac{1}{2 \pi i} \oint_{\Gamma} T(\lambda) d \lambda \in L\left(Z_{2}, Z_{1}\right)
$$

It holds trivially $S z \in D\left(A_{1}\right)=Z_{1}$ for any $z \in Z_{2}$. The same calculation, based on (6.2) but without the need of going through the Riemann sum as $D\left(A_{1}\right)=Z_{1}$, leads us to (6.3) which implies (6.4) for any $z \in D\left(A_{2}\right)$. The proof is complete.

In the next step, we remove the non-degeneracy assumption on $L$ and split the phase space $X$ into the direct sum of the hyperbolic (if any) and central subspaces of $J L$. In particular, the non-degeneracy of the quadratic form $\langle L \cdot, \cdot\rangle$ on the hyperbolic subspace $X_{u} \oplus X_{s}$ is of particular importance in the decomposition of $J L$.

Proposition 6.2. Assume (H1-3). There exist closed subspaces $X_{u, s, c} \subset X$ such that
(1) $X=X_{c} \oplus X_{u} \oplus X_{s}, X_{u, s} \subset D(J L), \operatorname{dim} X_{u}=\operatorname{dim} X_{s} \leq n^{-}(L)$, and $\operatorname{ker} L \subset X_{c}$;
(2) with respect to this decomposition, $J L$ and $L$ take the forms

$$
J L \longleftrightarrow\left(\begin{array}{ccc}
A_{c} & 0 & 0 \\
0 & A_{u} & 0 \\
0 & 0 & A_{s}
\end{array}\right), \quad L \longleftrightarrow\left(\begin{array}{ccc}
L_{X_{c}} & 0 & 0 \\
0 & 0 & B \\
0 & B^{*} & 0
\end{array}\right)
$$

(3) $B: X_{u} \rightarrow X_{s}^{*}$ is an isomorphism, $A_{c}$ is densely defined, closed, and the spectral sets satisfy $\sigma\left(A_{c}\right) \subset i \mathbf{R}$ and $\pm \operatorname{Re} \lambda>0$ for any $\lambda \in \sigma\left(A_{u, s}\right)$.

Proof. Let $X_{0}=\operatorname{ker} L$ and $Y=X_{-} \oplus X_{+}$where $X_{ \pm}$are given in Lemma 12.4. Let $P: X \rightarrow Y$ be the projection associated to $X=Y \oplus X_{0}$ and $J_{Y}=$ $P J P^{*}$. Lemma 12.3 implies that $\left(Y, L_{Y}, J_{Y}\right)$ satisfy assumptions (H1-3), with $L_{Y}$ being an isomorphism. Applying Proposition 6.1, we obtain closed subspaces $Y_{j}$, $j=1,2,3,4$, such that $X=X_{0} \oplus\left(\oplus_{j=1}^{4} Y_{j}\right)$ and $J L$ and $L$ take the forms

$$
J L \longleftrightarrow\left(\begin{array}{ccccc}
0 & \tilde{A}_{01} & \tilde{A}_{02} & \tilde{A}_{03} & \tilde{A}_{04} \\
0 & \tilde{A}_{1} & \tilde{A}_{12} & \tilde{A}_{13} & \tilde{A}_{14} \\
0 & 0 & \tilde{A}_{2} & 0 & \tilde{A}_{24} \\
0 & 0 & 0 & \tilde{A}_{3} & \tilde{A}_{34} \\
0 & 0 & 0 & 0 & \tilde{A}_{4}
\end{array}\right), \quad L \longleftrightarrow\left(\begin{array}{ccccc}
0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & \tilde{B} \\
0 & 0 & L_{Y_{2}} & 0 & 0 \\
0 & 0 & 0 & L_{Y_{3}} & 0 \\
0 & \tilde{B}^{*} & 0 & 0 & 0
\end{array}\right)
$$

where $\tilde{B}$ is an isomorphism, $L_{Y_{2}} \leq-\delta_{0}, L_{Y_{3}} \geq \delta_{0}$, for some $\delta_{0}>0$, and $\tilde{A}_{2,3}$ are anti-self-adjoint with respect to the equivalent inner products $\mp\left\langle L_{Y_{2,3}} \cdot, \cdot\right\rangle$ on $Y_{2,3}$. The upper triangular structure of $J L$ implies $\sigma(J L)=\{0\} \bigcup_{j=1}^{4} \sigma\left(\tilde{A}_{j}\right)$. Moreover, we have $\sigma\left(\tilde{A}_{2,3}\right) \subset i \mathbf{R}$ due to the anti-self-adjointness of $\tilde{A}_{2,3}$.

For $j=1,4$, as $\operatorname{dim} Y_{j}<\infty$, let $Y_{j}=Y_{j c} \oplus Y_{j h}$, where $Y_{j c}$ and $Y_{j h}$, are the eigenspaces of $\tilde{A}_{j}$ corresponding to all eigenvalues with zero and nonzero real parts, respectively. For any $x_{1,4} \in Y_{1,4}$, the above form of $J L$ and $L$ imply

$$
\begin{aligned}
\left\langle\tilde{B} x_{4}, \tilde{A}_{1} x_{1}\right\rangle+\left\langle\tilde{B} \tilde{A}_{4} x_{4}, x_{1}\right\rangle & =\left\langle L \tilde{A}_{1} x_{1}, x_{4}\right\rangle+\left\langle L \tilde{A}_{4} x_{4}, x_{1}\right\rangle \\
& =\left\langle L J L x_{1}, x_{4}\right\rangle+\left\langle L J L x_{4}, x_{1}\right\rangle=0
\end{aligned}
$$

Much as in the proof of Lemma 3.3, due to the difference in eigenvalues, we obtain

$$
\begin{equation*}
\left\langle\tilde{B} x_{4 h}, x_{1 c}\right\rangle=0=\left\langle\tilde{B} x_{4 c}, x_{1 h}\right\rangle, \forall x_{j c} \in Y_{j c}, x_{j h} \in Y_{j h}, j=1,4 \tag{6.5}
\end{equation*}
$$

Therefore, the non-degeneracy of $\tilde{B}$ implies that

$$
\begin{equation*}
\left\langle\tilde{B} x_{4 h}, x_{1 h}\right\rangle \text { and }\left\langle\tilde{B} x_{4 c}, x_{1 c}\right\rangle \text { are non-degenerate quadratic forms } \tag{6.6}
\end{equation*}
$$

on $Y_{1 h} \times Y_{4 h}$ and $Y_{1 c} \times Y_{4 c}$.
Applying Lemma 6.1 to $X_{0} \oplus Y_{1 h}$ and $\left.J L\right|_{X_{0} \oplus Y_{1 h}}$, we obtain a linear operator $S_{1}: Y_{1 h} \rightarrow X_{0}$ such that $X_{1 h} \triangleq\left(I+S_{1}\right) Y_{1 h} \subset D(J L)$ satisfies $J L\left(X_{1 h}\right)=X_{1 h}$ and $\sigma\left(\left.J L\right|_{X_{1 h}}\right)=\sigma\left(\left.\tilde{A}_{1}\right|_{Y_{1 h}}\right)$. Clearly, we still have the decomposition

$$
X=X_{0} \oplus Y_{1 c} \oplus X_{1 h} \oplus Y_{2} \oplus Y_{3} \oplus Y_{4 c} \oplus Y_{4 h}
$$

Applying again Lemma 6.1 to

$$
Z=X_{0} \oplus Y_{1 c} \oplus Y_{2} \oplus Y_{3} \oplus Y_{4 h}
$$

and the projection (with the kernel $X_{1 h} \oplus Y_{4 c}$ ) of $\left.J L\right|_{Z}$ to $Z$, we obtain a bounded linear operator

$$
S_{4}: Y_{4 h} \rightarrow X_{0} \oplus Y_{1 c} \oplus Y_{2} \oplus Y_{3}
$$

such that $X_{4 h} \triangleq\left(I+S_{4}\right) Y_{4 h} \subset D(J L)$ satisfies $J L\left(X_{4 h}\right) \subset X_{1 h} \oplus X_{4 h}$. Let $X_{h}=X_{1 h} \oplus X_{4 h}$, we have

$$
X_{h} \subset D(J L), J L\left(X_{h}\right)=X_{h}, \sigma\left(\left.J L\right|_{X_{h}}\right)=\sigma(J L) \backslash i \mathbf{R}=\sigma\left(\left.\tilde{A}_{1}\right|_{Y_{1 h}}\right) \cup \sigma\left(\left.\tilde{A}_{4}\right|_{Y_{4 h}}\right)
$$

According to (6.5) and the form of $L$, it holds

$$
\left\langle L\left(I+S_{1}\right) x_{1 h},\left(I+S_{4}\right) x_{4 h}\right\rangle=\left\langle L x_{1 h}, x_{4 h}\right\rangle=\left\langle\tilde{B} x_{4 h}, x_{1 h}\right\rangle .
$$

Therefore, we obtain the non-degeneracy of $\langle L \cdot, \cdot\rangle$ on $X_{h}$ from (6.6) and the construction of $X_{h}$. Let $X_{h}=X_{u} \oplus X_{s}$, where $X_{u, s}$ are the eigenspaces of all eigenvalues $\lambda \in \sigma\left(\left.J L\right|_{X_{h}}\right)$ with $\pm \operatorname{Re} \lambda>0$. Lemma 3.3 implies $\langle L u, v\rangle=0$ on both $X_{u}$ and $X_{s}$ and thus $\langle L \cdot, \cdot\rangle$ is a non-degenerate quadratic form on $X_{u} \times X_{s}$ due to the non-degeneracy of $\langle L \cdot, \cdot\rangle$ on $X_{h}$. This also yields

$$
\operatorname{dim} X_{u}=\operatorname{dim} X_{s}=\frac{1}{2} \operatorname{dim} X_{h} \leq \operatorname{dim} Y_{1} \leq n^{-}(L)
$$

Let

$$
X_{c}=X_{h}^{\perp_{L}}=\left\{u \in X \mid\langle L u, v\rangle=0, \forall v \in X_{h}\right\} .
$$

By Lemmas 12.2 and 3.2, $X=X_{h} \oplus X_{c}$ and $X_{c}$ is invariant under $e^{t J L}$. Therefore, $J L$ is densely defined on $X_{c}$ and $A_{c}\left(D\left(A_{c}\right) \cap X_{c}\right) \subset X_{c}$ where $A_{c}=\left.J L\right|_{X_{c}}$. Moreover, from the non-degeneracy of $\left\langle L^{\cdot}, \cdot\right\rangle$ on $X_{h}$, it is straightforward to show that $X_{c}$ can be written as a graph of a bounded linear operator from $X_{0} \oplus Y_{1 c} \oplus Y_{2} \oplus Y_{3} \oplus Y_{4 c}$ to $X_{h}$. Therefore, due to the upper triangular structure of $J L$, the spectrum $\sigma\left(\left.J L\right|_{X_{c}}\right)$ is given by the union of the spectrum of those diagonal blocks of $J L$ complementary to $Y_{1 h}$ and $Y_{4 h}$ and thus $\sigma\left(\left.J L\right|_{X_{c}}\right) \subset i \mathbf{R}$.

As a by-product, we prove the symmetry of eigenvalues of $\sigma(J L)$.
Corollary 6.1. Suppose $\lambda \in \sigma(J L)$.
(i) If $\lambda \in \sigma(J L) \backslash i \mathbf{R}$, then $\lambda$ is an isolated eigenvalue of finite algebraic multiplicity. Its eigenspace consists of generalized eigenvectors only. Moreover, let $m_{\lambda}$ to be the algebraic multiplicity of $\lambda$, then

$$
\begin{equation*}
n^{-}\left(\left.L\right|_{E_{\lambda} \oplus E_{-\bar{\lambda}}}\right)=\operatorname{dim}\left(E_{\lambda}\right)=m_{\lambda} . \tag{6.7}
\end{equation*}
$$

(ii) If $\lambda$ is an eigenvalues of $J L$, then $\pm \lambda, \pm \bar{\lambda}$ are also eigenvalues of $J L$. Moreover, for any integer $k>0$, $\operatorname{dim} \operatorname{ker}(J L-a)^{k}$ are the same for $a= \pm \lambda, \pm \bar{\lambda}$.

For an eigenvalue $\lambda \in i \mathbf{R}$, it may happen $\operatorname{dim} \operatorname{ker}(J L-\lambda)=\infty$.
Proof. According to Lemma 3.6, we only need to prove $\lambda \in \sigma(J L) \backslash i \mathbf{R}$ implies that $\lambda$ is an isolated eigenvalue of finite multiplicity and $\operatorname{dim} \operatorname{ker}(J L-\lambda)^{k}=$ $\operatorname{dim} \operatorname{ker}(J L+\bar{\lambda})^{k}$.

In fact, if $\lambda \in \sigma(J L) \backslash i \mathbf{R}$, then Proposition 6.2 implies that $\lambda \in \sigma\left(A_{u}\right) \cup \sigma\left(A_{s}\right)$. As $A_{u, s}$ are finite dimensional matrices, $\lambda$ must be an isolated eigenvalue of $J L$ with finite algebraic multiplicity. Moreover, from the blockwise forms of $L$ and $J L$ and $J^{*}=-J$, it is easy to compute

$$
J \longleftrightarrow\left(\begin{array}{ccc}
J_{X_{c}} & 0 & 0 \\
0 & 0 & A_{u}\left(B^{*}\right)^{-1} \\
0 & A_{s} B^{-1} & 0
\end{array}\right)
$$

Again since $J^{*}=-J$, we have $A_{s}=-B^{-1} A_{u}^{*} B$. As $A_{u, s}$ are finite dimensional matrices and eigenvalues of $J L$ with positive (or negative) real parts coincide with eigenvalues of $A_{u}$ (or $A_{s}$ ), the statement in the corollary follows from this similarity immediately.

Since by Proposition 6.2 $\left.L\right|_{X_{u} \oplus X_{s}}$ is non-degenerate, formula (6.7) follows from Lemma 4.2 in the finite dimensional case.

Proof of Theorem 2.1. Let $X_{5,6}=X_{u, s}$ and $J_{X_{c}}=P_{c} J P_{c}^{*}$, where $X_{u, s, c}$ are obtained in Proposition 6.2 and $P_{c}: X \rightarrow X_{c}$ be the projection associated to $X=X_{c} \oplus X_{u} \oplus X_{s}$. According to Lemma 12.3, ( $\left.X_{c}, L_{X_{c}}, J_{X_{c}}\right)$ satisfy assumption (H1-3) as well. Since Proposition 6.2 also ensures the non-degeneracy of $L_{X_{5} \oplus X_{6}}$ and $\operatorname{dim} X_{5,6} \leq n^{-}(L)$, the finite dimensional results in Chapter 4 (Lemma 3.6 and 4.2) imply the symmetry between the spectra $\sigma\left(A_{5}\right)$ and $\sigma\left(A_{6}\right)$ and $n^{-}\left(\left.L\right|_{X_{5} \oplus X_{6}}\right)=$ $\operatorname{dim} X_{5}$. Therefore, we obtain, from the $L$-orthogonality between $X_{c}$ and $X_{u} \oplus X_{s}$,

$$
n^{-}\left(L_{X_{c}}\right)=n^{-}(L)-\operatorname{dim} X_{5} .
$$

Recall $X_{0}=\operatorname{ker} L=\operatorname{ker} L_{X_{c}} \subset X_{c}$. Let $X_{ \pm}$be given by Lemma 12.4 applied to $\left(X_{c}, L_{X_{c}}, J_{X_{c}}\right), Y=X_{+} \oplus X_{-}, P_{Y}: X \rightarrow Y$ be the associated projection, and $J_{Y}=P_{Y} J P_{Y}^{*}$. Again Lemma 12.3 implies $\left(Y, L_{Y}, J_{Y}\right)$ satisfy $(\mathbf{H 1} \mathbf{- 3})$ with $L_{Y}$ being an isomorphism. Applying Proposition 6.1 to $Y$ and we obtain subspaces $\tilde{X}_{j}$, $j=1,2,3,4$. To ensure the orthogonality between $X_{0}=\operatorname{ker} L$ and $X_{j}, j=1,2,3,4$, we modify the definition of $X_{j}$ as

$$
X_{j}=\left\{u \in X_{0} \oplus \tilde{X}_{j} \mid(u, v)=0, \forall v \in \operatorname{ker} L\right\}, \quad j=1,2,3,4
$$

It is straightforward to verify the desired properties of the decomposition $X=$ $\oplus_{j=0}^{6} X_{j}$ by using Propositions 6.1 and 6.2. The proof of Theorem 2.1 is complete.

To finish this chapter, we give the following lemma on the $L$-orthogonality between certain eigenspaces defined by spectral integrals.

LEMMA 6.2. Let $\Omega \subset \mathbf{C}$ be an open subset symmetric about $i \mathbf{R}$ with smooth boundary $\Gamma=\partial \Omega$ and compact closure such that $\Gamma \cap \sigma(J L)=\emptyset$. Let

$$
P=\frac{1}{2 \pi i} \oint_{\Gamma}(z-J L)^{-1} d z .
$$

and then it holds that $\langle L(I-P) u, P v\rangle=0$, for any $u, v \in X$.
The above $P$ is simply the standard spectral projection operator.
Proof. We first observe for any $w, w^{\prime} \in X,(12.8)$ and (12.10) imply

$$
\begin{align*}
& \frac{1}{2 \pi i} \oint_{\Gamma}\left\langle L w,(z-J L)^{-1} w^{\prime}\right\rangle d z=\left\langle L w, P w^{\prime}\right\rangle  \tag{6.8}\\
& \frac{1}{2 \pi i} \oint_{\Gamma}\left\langle L(z-J L)^{-1} w, w^{\prime}\right\rangle d \bar{z}=-\left\langle L P w, w^{\prime}\right\rangle \tag{6.9}
\end{align*}
$$

where the first equality is used in the derivation of the second equality. Here the $d \bar{z}$ and the minus sign in the second equality are due to the anti-linear nature of $L$ in (12.10).

Let $\Omega_{1} \subset \Omega$ be an open subset symmetric about $i \mathbf{R}$ such that $\Gamma_{1}=\partial \Omega_{1} \subset \Omega$ is smooth and $\sigma(J L) \cap\left(\Omega \backslash \Omega_{1}\right)=\emptyset$. Clearly,

$$
P=\frac{1}{2 \pi i} \oint_{\Gamma_{1}}(z-J L)^{-1} d z
$$

due to the analyticity of $(z-J L)^{-1}$. Denote

$$
\tilde{u}(z)=(z-J L)^{-1} u, \tilde{v}(z)=(z-J L)^{-1} v, \quad \forall z \notin \sigma(J L) .
$$

For $z_{1}, z_{2} \notin \sigma(J L)$ satisfying $\bar{z}_{1}+z_{2} \neq 0$, one may compute using (12.8) and (12.10)

$$
\begin{aligned}
& \frac{1}{\bar{z}_{1}+z_{2}}\left(\left\langle L\left(z_{1}-J L\right)^{-1} u, v\right\rangle+\left\langle L u,\left(z_{2}-J L\right)^{-1} v\right\rangle\right) \\
= & \frac{1}{\bar{z}_{1}+z_{2}}\left(\left\langle L \tilde{u}\left(z_{1}\right),\left(z_{2}-J L\right) \tilde{v}\left(z_{2}\right)\right\rangle+\left\langle L\left(z_{1}-J L\right) \tilde{u}\left(z_{1}\right), \tilde{v}\left(z_{2}\right)\right\rangle\right) \\
= & \left\langle L \tilde{u}\left(z_{1}\right), \tilde{v}\left(z_{2}\right)\right\rangle=\left\langle L\left(z_{1}-J L\right)^{-1} u,\left(z_{2}-J L\right)^{-1} v\right\rangle .
\end{aligned}
$$

Due to the definition of $\Gamma_{1}$ and its symmetry about the imaginary axis, $\bar{z}_{1}+z_{2} \neq 0$ for any $z_{1} \in \Gamma$ and $z_{2} \in \Gamma_{1}$. Integrating the above equality along these curves, where $\Gamma_{1}$ is enclosed in $\Gamma$, we obtain from the Cauchy integral theorem and (6.8) and (6.9)

$$
\begin{aligned}
\langle L P u, P v\rangle= & \frac{-1}{(2 \pi i)^{2}} \oint_{\Gamma} \oint_{\Gamma_{1}}\left\langle L\left(z_{1}-J L\right)^{-1} u,\left(z_{2}-J L\right)^{-1} v\right\rangle d z_{2} d \bar{z}_{1} \\
= & \frac{-1}{(2 \pi i)^{2}} \oint_{\Gamma} \oint_{\Gamma_{1}} \frac{1}{\bar{z}_{1}+z_{2}}\left\langle L\left(z_{1}-J L\right)^{-1} u, v\right\rangle d z_{2} d \bar{z}_{1} \\
& +\frac{-1}{(2 \pi i)^{2}} \oint_{\Gamma_{1}} \oint_{\Gamma} \frac{1}{\bar{z}_{1}+z_{2}}\left\langle L u,\left(z_{2}-J L\right)^{-1} v\right\rangle d \bar{z}_{1} d z_{2} .
\end{aligned}
$$

Since $-\bar{z}_{1}$ is not enclosed in $\Gamma_{1}$ while $-\bar{z}_{1}$ is enclosed in $\Gamma$, the above first integral vanishes and the we obtain from (6.8) and the Cauchy integral theorem

$$
\langle L P u, P v\rangle=\langle L u, P v\rangle
$$

This proves the lemma.
The above lemma implies that $\langle L u, v\rangle=0$ for any $u \in \operatorname{ker} P$ and $v \in P X$, where $X=P X \oplus$ ker $P$ is a spectral decomposition of $X$ invariant under $J L$. As a corollary, we give the following extension of Lemma 3.3.

Let $\tilde{\sigma} \subset \sigma(J L)$ be compact and also open in the relative topology of $\sigma(J L)$, namely $\tilde{\sigma}$ is isolated in $\sigma(J L)$. There exists an open domain $\Omega \subset \mathbf{C}$ with compact closure and smooth boundary such that $\Omega \cap \sigma(J L)=\tilde{\sigma}$. Let

$$
P_{\tilde{\sigma}}=\frac{1}{2 \pi i} \oint_{\partial \Omega}(z-J L)^{-1} d z, \quad X_{\tilde{\sigma}}=P_{\tilde{\sigma}} X, \quad X_{\tilde{\sigma}^{c}}=\operatorname{ker} P_{\tilde{\sigma}}
$$

According to the Cauchy integral theorem, the projection operator $P_{\tilde{\sigma}}$ as well as the above subspaces, which are invariant under $J L$, are independent of the choice of $\Omega$ and $J L P_{\tilde{\sigma}}=P_{\tilde{\sigma}} J L$. Moreover,

$$
\sigma\left(\left.J L\right|_{X_{\tilde{\sigma}}}\right)=\tilde{\sigma}, \quad \sigma\left(\left.J L\right|_{X_{\tilde{\sigma}^{c}}}\right)=\sigma(J L) \backslash \tilde{\sigma} .
$$

Corollary 6.2. Suppose $\sigma_{j} \subset \sigma(J L), j=1,2$, are compact and also open in the relative topology of $\sigma(J L)$. In addition, assume

$$
\sigma_{1} \cap \tilde{\sigma}_{2}=\emptyset, \text { where } \tilde{\sigma}_{2}=\left\{\lambda \in \mathbf{C} \mid \lambda \in \sigma_{2} \text { or } \bar{\lambda} \in \sigma_{2}\right\} .
$$

Then $\langle L u, v\rangle=0$ for any $u \in X_{\sigma_{1}}$ and $v \in X_{\sigma_{2}}$ where $X_{\sigma_{1,2}}$ are defined as in the above.

Proof. According to our assumptions, there exists an open domain $\Omega \subset \mathbf{C}$, symmetric about $i \mathbf{R}$ with smooth boundary and compact closure such that $\Omega \cap$ $\sigma(J L)=\tilde{\sigma}_{2}$ and $\partial \Omega \cap \sigma(J L)=\emptyset$. The corollary follows from Lemma 6.2 and the facts $X_{\sigma_{1}} \subset \operatorname{ker} P_{\tilde{\sigma}_{2}}$ and $X_{\sigma_{2}} \subset P_{\tilde{\sigma}_{2}} X$.

## CHAPTER 7

## Exponential trichotomy

We prove Theorem 2.2 on the exponential trichotomy in this chapter. The proof is based on the decomposition Theorem 2.1 and we follow the notations there.

Let

$$
E^{u}=X_{5}, \quad E^{s}=X_{6}, \quad E^{c}=\oplus_{j=0}^{4} X_{j}
$$

where $X_{j}, j=0, \ldots, 6$, are given by Theorem 2.1. Based on Theorem 2.1, it only remains to prove the growth estimates.

Since $A_{2,3}$ are anti-self-adjoint with respect to the equivalent inner product $\mp\left\langle L_{X_{2,3}} \cdot, \cdot\right\rangle$, there exists a constant $C>0$ such that

$$
\begin{equation*}
\left|e^{t A_{2}}\right|,\left|e^{t A_{3}}\right| \leq C, \forall t \in \mathbf{R} \tag{7.1}
\end{equation*}
$$

Since $\operatorname{dim} X_{5}=\operatorname{dim} X_{6}<\infty$ and $\sigma\left(A_{5}\right)=-\sigma\left(A_{6}\right)$, it is clear

$$
\begin{align*}
& \left|e^{t A_{5}}\right| \leq C\left(1+|t|^{\operatorname{dim} X_{5}-1}\right) e^{\lambda_{u} t}, \forall t<0 \\
& \left|e^{t A_{6}}\right| \leq C\left(1+|t|^{\operatorname{dim} X_{6}-1}\right) e^{-\lambda_{u} t}, \forall t>0 \tag{7.2}
\end{align*}
$$

for some $C>0$ and $\lambda_{u}=\min \left\{\operatorname{Re} \lambda \mid \lambda \in \sigma\left(A_{5}\right)\right\}$. Finally, as $\operatorname{dim} X_{1}=\operatorname{dim} X_{4}<\infty$ and $\sigma\left(A_{1,4}\right) \subset i \mathbf{R}$, we also have

$$
\begin{equation*}
\left|e^{t A_{1,4}}\right| \leq C\left(1+|t|^{\operatorname{dim} X_{1}-1}\right), \forall t \in \mathbf{R} . \tag{7.3}
\end{equation*}
$$

For any $x \in X$, write

$$
e^{t J L} x=\sum_{j=0}^{6} x_{j}(t), \quad x_{j}(t) \in X_{j}
$$

where $X_{j}, j=0, \ldots, 6$, are given by Theorem 2.1. One can write down the equations explicitly:

$$
\left\{\begin{array}{l}
\partial_{t} x_{0}=A_{01} x_{1}+A_{02} x_{2}+A_{03} x_{3}+A_{04} x_{4}  \tag{7.4}\\
\partial_{t} x_{1}=A_{1} x_{1}+A_{12} x_{2}+A_{13} x_{3}+A_{14} x_{4} \\
\partial_{t} x_{2}=A_{2} x_{2}+A_{24} x_{4} \\
\partial_{t} x_{3}=A_{3} x_{3}+A_{34} x_{4} \\
x_{j}(t)=e^{t A_{j}} x_{j}(0), j=4,5,6 .
\end{array}\right.
$$

For $j=2,3$, we obtain from Theorem 2.1 and inequalities (7.1) and (7.3) that

$$
\begin{align*}
\left\|x_{j}(t)\right\| & =\left\|e^{t A_{j}} x_{j}(0)+\int_{0}^{t} e^{(t-\tau) A_{j}} A_{j 4} e^{\tau A_{4}} x_{4}(0) d \tau\right\|  \tag{7.5}\\
& \leq C\left(\left\|x_{j}(0)\right\|+\left(1+|t|^{\operatorname{dim} X_{1}}\right)\left\|x_{4}(0)\right\|\right)
\end{align*}
$$

for some $C>0$. Regrading $x_{1}(t)$, we have from (7.1), (7.3), and (7.5)

$$
\begin{align*}
\left\|x_{1}(t)\right\| & \leq\left\|e^{t A_{1}} x_{1}(0)+\int_{0}^{t} e^{(t-\tau) A_{1}}\left(A_{12} x_{2}(\tau)+A_{13} x_{3}(\tau)+A_{14} e^{\tau A_{4}} x_{4}(0)\right) d \tau\right\| \\
& \leq C\left(1+|t|^{\operatorname{dim} X_{1}-1}+\int_{0}^{|t|} 1+|t-\tau|^{\operatorname{dim} X_{1}-1}|\tau|^{\operatorname{dim} X_{1}} d \tau\right)\|x(0)\| \\
& \leq C\left(1+|t|^{2 \operatorname{dim} X_{1}}\right)\|x(0)\| . \tag{7.6}
\end{align*}
$$

Much as on the above we also have

$$
\begin{equation*}
\left\|x_{0}(t)\right\| \leq C\left(1+|t|^{2 \operatorname{dim} X_{1}+1}\right)\|x(0)\| \tag{7.7}
\end{equation*}
$$

The above inequalities prove the desired exponential trichotomy estimates.
Finally, repeatedly applying $J L$ to equation (2.1) and using the above inequalities yield the trichotomy estimates in the graph norms on $D\left((J L)^{k}\right)$.

## CHAPTER 8

## The index theorems and the structure of $E_{i \mu}$

Our goal in this chapter is to complete the proof of the index theorems and related properties.

### 8.1. Proof of Theorem 2.3: the index counting formula

The symmetry of $\sigma(J L)$, the eigenvalues of $J L$, and the dimensions of the spaces of generalized eigenvectors have been proved in Lemma 3.6 and Corollary 6.1. The index formula (2.13) will be proved in the next two lemmas. Recall the notations $n^{-}\left(\left.L\right|_{Y}\right)$ and $n^{\leq 0}\left(\left.L\right|_{Y}\right)$ for a subspace $Y \subset X$ and indices $k_{r}, k_{c}, k^{\leq 0}(i \mu)$, $k_{i}^{\leq 0}, k_{0}^{\leq 0}$ etc. defined in Section 2.4.

Lemma 8.1. Under hypotheses (H1-3), it holds

$$
k_{r}+2 k_{c}+2 k_{i}^{\leq 0}+k_{0}^{\leq 0} \geq n^{-}(L)
$$

Proof. Let $X_{j}, j=0, \ldots, 6$, be the closed subspaces constructed in Theorem 2.1 and $Z=\oplus_{j=0}^{2} X_{j}$. From Theorem 2.1, $Z$ is an invariant subspace of $J L$ containing ker $L$ satisfying $\sigma\left(\left.J L\right|_{Z}\right) \subset i \mathbf{R}$. For any eigenvalue $i \mu \in \sigma\left(\left.J L\right|_{Z}\right)$, let $E_{i \mu}(Z)=E_{i \mu} \cap Z$ be the subspace of generalized eigenvectors of $i \mu$ in $Z$, and denote the corresponding non-positive index of $\left.L\right|_{Z}$ by

$$
k_{i}^{\leq 0}(Z)=\Sigma_{i \mu \in \sigma\left(\left.J L\right|_{Z}\right) \cap i \mathbf{R}^{+}} k^{\leq 0}(i \mu, Z)
$$

where $k^{\leq 0}(i \mu, Z)=n^{\leq 0}\left(\left.L\right|_{E_{i \mu}(Z)}\right)$.
On the one hand, for any eigenvalue $i \mu \neq 0$, it clearly holds $E_{i \mu}(Z) \subset E_{i \mu}$ and thus $k^{\leq 0}(i \mu, Z) \leq k^{\leq 0}(i \mu)$. Therefore, we have $k_{i}^{\leq 0}(Z) \leq k_{i}^{\leq 0}$. For the same reason, we also have $k_{0}^{\leq 0}(Z) \leq k_{0}^{\leq 0}$ as $\operatorname{ker} L \subset E_{0}(Z)$, where $k_{0}^{\leq 0}(Z)$ has a similar definition as $k_{0}^{\leq 0}$ (defined in (2.12)) except applied to $E_{0}(Z)$ instead of $E_{0}$. From Theorem 2.1 and the finite dimensionality of $X_{5}$, it is clear $k_{r}+2 k_{c}=\operatorname{dim} X_{5}$. Consequently, we obtain

$$
\begin{equation*}
k_{r}+2 k_{c}+2 k_{i}^{\leq 0}+k_{0}^{\leq 0} \geq \operatorname{dim} X_{5}+2 k_{i}^{\leq 0}(Z)+k_{0}^{\leq 0}(Z) \tag{8.1}
\end{equation*}
$$

On the other hand, due to the finite dimensionality of $X_{j}, j=1,2$, and the blockwise upper triangular form of $J L$, we have $Z=\oplus_{i \mu \in \sigma\left(\left.J L\right|_{Z}\right) \cap i \mathbf{R}} E_{i \mu}(Z)$. Moreover, since $L$ is non-positive on $Z$ according to Theorem 2.1, we have

$$
2 k_{i}^{\leq 0}(Z)+k_{0}^{\leq 0}(Z)=\operatorname{dim} X_{1}+\operatorname{dim} X_{2}=n^{-}(L)-\operatorname{dim} X_{5}
$$

Combining it with (8.1), we obtain the conclusion of the lemma.
Lemma 8.2. Under hypotheses (H1-3), it holds

$$
k_{r}+2 k_{c}+2 k_{i}^{\leq 0}+k_{0}^{\leq 0} \leq n^{-}(L) .
$$

Proof. Let $X_{j}, j=0, \ldots, 6$, be the closed subspaces constructed in Theorem 2.1 and $Y=\oplus_{j=1}^{6} X_{j}$. Let $P_{Y}$ be the projection associated to $X=\operatorname{ker} L \oplus$ $Y$. Lemma 12.3 implies that $\left(Y, L_{Y}, J_{Y}\right)$ satisfies assumptions (H1-3), where $n^{-}\left(L_{Y}\right)=n^{-}(L)$. The definitions of $J_{Y}$ and $L_{Y}$ also imply $J_{Y} L_{Y}=P_{Y}(J L)$.

Let $i \mu \in \sigma(J L) \cap i \mathbf{R}^{+}$. By the definition of $k^{\leq 0}(i \mu)$, there exists a subspace $E_{i \mu}^{\leq 0} \subset E_{i \mu}$ such that $\operatorname{dim} E_{i \mu}^{\leq 0}=k^{\leq 0}(i \mu)$ and $\langle L u, u\rangle \leq 0$, for all $u \in E_{i \mu}^{\leq 0}$. Since $\mu \neq 0$ and thus $E_{i \mu} \cap \operatorname{ker} L=\{0\}$, we have $\operatorname{dim} P_{Y} E_{i \mu}^{\leq 0}=\operatorname{dim} E_{i \mu}^{\leq 0}$. For $\mu<0$, let $E_{i \mu}^{\leq 0}=\left\{\bar{u} \mid u \in E_{-i \mu}^{\leq 0}\right\}$. For $\mu=0$, let $\tilde{E}_{0}=E_{0} \cap Y$ where clearly $E_{0}=\operatorname{ker} L \oplus \tilde{E}_{0}$. There exists a subspace $E_{0}^{\leq 0} \subset \tilde{E}_{0}$ such that $\operatorname{dim} E_{0}^{\leq 0}=k_{0}^{\leq 0}$ and $\langle L u, u\rangle \leq 0$, for all $u \in E_{0}^{\leq 0}$. Let

$$
W=X_{5} \oplus E_{0}^{\leq 0} \oplus\left(\oplus_{i \mu \in \sigma J L \cap i \mathbf{R}} P_{Y} E_{i \mu}^{\leq 0}\right) \subset Y
$$

It is clearly (the complexification of) a real subspace of $Y$ satisfying $\bar{u} \in W$ for all $u \in W$. Theorem 2.1 implies

$$
\operatorname{dim} W=\operatorname{dim} X_{5}+k_{0}^{\leq 0}+2 k_{i}^{\leq 0}=k_{r}+2 k_{c}+k_{0}^{\leq 0}+2 k_{i}^{\leq 0}
$$

From Lemma 3.3, we have $X_{5}$ and $P_{Y} E_{i \mu}^{\leq 0}(i \mu \in \sigma J L \cap i \mathbf{R})$ are mutually $L$-orthogonal. Therefore, our construction of $W$ yields that $\langle L u, u\rangle \leq 0$ for all $u \in W \subset Y$. Applying Theorem 5.1 to $\left(Y, L_{Y}, J_{Y}\right)$ implies $\operatorname{dim} W \leq n^{-}\left(L_{Y}\right)=n^{-}(L)$ and thus the lemma is proved.

### 8.2. Structures of subspaces $E_{i \mu}$ of generalized eigenvectors

In this section, we will prove Propositions 2.1 and 2.2. We complete the proof in several steps.

Lemma 8.3. Let $i \mu \in \sigma(J L) \cap i \mathbf{R}$ and $E \subset E_{i \mu}$ be a closed subspace such that $J L(E) \subset E$. In addition to (H1-3), assume $\langle L \cdot, \cdot\rangle$ is non-degenerate (in the sense of (2.4)) on both $X$ and $E$. Then there exist closed subspaces $E^{1}, \tilde{E} \subset E$ such that $E=E^{1} \oplus \tilde{E}$ and $L, J L$ take the following forms on $E$

$$
\langle L \cdot, \cdot\rangle \longleftrightarrow\left(\begin{array}{cc}
L_{E^{1}} & 0 \\
0 & L_{\tilde{E}}
\end{array}\right), \quad J L \longleftrightarrow\left(\begin{array}{cc}
i \mu & 0 \\
0 & \tilde{A}
\end{array}\right)
$$

and $\operatorname{ker}(J L-i \mu) \cap \tilde{E} \subset(J L-i \mu) \tilde{E}$ with non-degenerate $L_{E^{1}}$ and $L_{\tilde{E}}$ and

$$
\operatorname{dim} \tilde{E} \leq 3\left(n^{-}\left(\left.L\right|_{E}\right)-n^{-}\left(\left.L\right|_{E^{1}}\right)\right), \operatorname{dim}((J L-i \mu) E) \leq 2\left(n^{-}\left(\left.L\right|_{E}\right)-n^{-}\left(\left.L\right|_{E^{1}}\right)\right)
$$

REMARK 8.1. The property $\operatorname{ker}(J L-i \mu) \cap \tilde{E} \subset(J L-i \mu) \tilde{E}$, or equivalently $\operatorname{ker}(\tilde{A}-i \mu) \subset(\tilde{A}-i \mu) \tilde{E}$, is equivalent to that the Jordan canonical form of $\tilde{A}$ contains only nontrivial Jordan blocks.

Proof. From Lemma 3.5, $E_{i \mu}=\operatorname{ker}(J L-i \mu)^{K}$ for some $K>0$ and $J L$ : $E_{i \mu} \rightarrow E_{i \mu}$ is a bounded operator. Let

$$
\begin{aligned}
& E^{0}=\{u \in E \cap \operatorname{ker}(J L-i \mu) \mid\langle L u, v\rangle=0, \forall v \in E \cap \operatorname{ker}(J L-i \mu)\} \\
& E^{1}=\left\{u \in E \cap \operatorname{ker}(J L-i \mu) \mid(u, v)=0, \forall v \in E^{0}\right\}
\end{aligned}
$$

Obviously, $\operatorname{ker}(J L-i \mu) \cap E=E^{0} \oplus E^{1}$. Moreover, for any $u \in E^{1} \backslash\{0\}$, there must exist $v \in E^{1}$ such that $\langle L u, v\rangle \neq 0$, otherwise it would lead to $u \in E^{0}$, a contradiction. Applying statement 2 of Lemma 12.2 to $Y=E^{1}$, we obtain that
$\langle L \cdot, \cdot\rangle$ is non-degenerate on $E^{1}$. Since $\langle L \cdot, \cdot\rangle$ is assumed to be non-degenerate on both $X$ and $E$, we apply statement 1 of Lemma 12.2 to obtain

$$
\begin{equation*}
X=E^{1} \oplus\left(E^{1}\right)^{\perp_{L}} \text { and } E=E^{1} \oplus \tilde{E}, \text { where } \tilde{E}=E \cap\left(E^{1}\right)^{\perp_{L}} \tag{8.2}
\end{equation*}
$$

Here $\left(E^{1}\right)^{\perp_{L}} \subset X$ is the subspace $L$-perpendicular to $E^{1}$. Clearly, $E^{0} \subset \tilde{E}$.
Claim. 1.) $\operatorname{dim} \tilde{E}<\infty$, 2.) $\langle L \cdot, \cdot\rangle$ is non-degenerate on $\tilde{E}$, and 3.) $J L(\tilde{E}) \subset$ $\tilde{E}$.

The invariance of $\tilde{E}$ under $J L$ follows directly from the invariance of $E$ and $E^{1}$ and Lemma 3.2. The non-degeneracy of $\langle L \cdot, \cdot\rangle$ on both $E$ and $E^{1}$ implies that $\langle L \cdot, \cdot\rangle$ is non-degenerate on $\tilde{E}$ as well. To complete the proof of the claim, we only need to prove $\operatorname{dim} \tilde{E}<\infty$.

On the one hand, from the above definitions, $\langle L u, v\rangle=0$ for any $u, v \in E^{0}$. The non-degeneracy of $\langle L \cdot, \cdot\rangle$ on $E$ and Theorem 5.1 along with Remark 5.1 imply

$$
\begin{equation*}
\operatorname{dim} E^{0} \leq n^{-}\left(\left.L\right|_{\tilde{E}}\right) \tag{8.3}
\end{equation*}
$$

On the other hand, it is clear from the definitions of $\tilde{E}$ and the non-degeneracy of $\langle L \cdot, \cdot\rangle$ on $E^{1}$ that

$$
\begin{equation*}
\tilde{E} \cap \operatorname{ker}(J L-i \mu)=E^{0} \tag{8.4}
\end{equation*}
$$

Moreover, from Lemma 3.5, $E \subset E_{i \mu}=\operatorname{ker}(J L-i \mu)^{K}$ for some $K>0$, and each Jordan chain in $\tilde{E}$ contains a vector in $E^{0}$, we obtain

$$
\operatorname{dim} \tilde{E} \leq K \operatorname{dim} E^{0} \leq K n^{-}\left(\left.L\right|_{E}\right)
$$

from the invariance of $\tilde{E}$ under $J L$. The claim is proved.
Now we complete the proof of the lemma by reducing it to a finite dimensional problem satisfying our framework. Firstly, to replace $\tilde{E}$ by the complexification of some real Hilbert space, let

$$
\tilde{E}^{R}=\{u+\bar{v} \mid u, v \in \tilde{E}\}
$$

which satisfies $\bar{u} \in \tilde{E}^{R}$ for any $u \in \tilde{E}^{R}$. Since $u \in E_{i \mu}$ implies $\bar{u} \in E_{-i \mu}$, we have $\tilde{E}^{R}=\tilde{E}$ if $\mu=0$. If $\mu \neq 0$, from (12.12) and Lemma 3.3 we obtain

$$
\langle L \bar{u}, v\rangle=0,\langle L \bar{u}, \bar{v}\rangle=\overline{\langle L u, v\rangle}, \quad \forall u, v \in \tilde{E}
$$

Therefore, $\tilde{E}^{R}$ satisfies the same properties as in the above claim whether $\mu=0$ or not. Using the non-degeneracy of $\langle L \cdot, \cdot\rangle$ on $X$ and $\tilde{E}^{R}$, and applying Lemma 12.3 to the splitting $X=\tilde{E}^{R} \oplus\left(\tilde{E}^{R}\right)^{\perp_{L}}$ with the associated projections $P_{\tilde{E}^{R}}$ and $I-P_{\tilde{E}^{R}}$, we have that the combination $\left(\tilde{E}^{R}, L_{\tilde{E}^{R}}, J_{\tilde{E}^{R}}\right)$ satisfies assumptions $(\mathbf{H 1} \mathbf{- 3})$, where $J_{\tilde{E}^{R}}=P_{\tilde{E}^{R}} J P_{\tilde{E}^{R}}^{*}$. We may apply Proposition 2.2 , whose finite dimensional case under the non-degeneracy assumption on $\langle L \cdot, \cdot\rangle$ has been proved in Chapter 4. As $\tilde{E}=\tilde{E}^{R} \cap \operatorname{ker}(J L-i \mu)^{K}$, that canonical form implies

$$
\operatorname{dim}((J L-i \mu) \tilde{E}) \leq 2 n^{-}\left(\left.L\right|_{\tilde{E}}\right), \quad \operatorname{ker}(J L-i \mu) \cap \tilde{E} \subset(J L-i \mu) \tilde{E}
$$

where (8.4) is also used along with the canonical form. We notice $(J L-i \mu) E=$ $(J L-i \mu) \tilde{E}$, as $E^{1} \subset \operatorname{ker}(J L-i \mu)$, and thus

$$
\operatorname{dim}((J L-i \mu) E) \leq 2 n^{-}\left(\left.L\right|_{\tilde{E}}\right)=2\left(n^{-}\left(\left.L\right|_{E}\right)-n^{-}\left(L_{E^{1}}\right)\right)
$$

The block forms of $L$ and $J L$ follow from the $L$-orthogonality and the invariance of the splitting $E=E^{1} \oplus \tilde{E}$. Finally, the estimate on $\operatorname{dim} \tilde{E}$ follows from the above inequality and (8.3) and (8.4). The proof is complete.

Next we study $E_{i \mu}$ by assuming the non-degeneracy of $L$.
Lemma 8.4. In addition to (H1-3), assume $\langle L \cdot, \cdot\rangle$ is non-degenerate. Let $i \mu \in$ $\sigma(J L) \cap i \mathbf{R}$. There exist subspaces $E^{D, 1, G} \subset E_{i \mu}$ such that

$$
\begin{gathered}
E_{i \mu}=E^{D} \oplus E^{1} \oplus E^{G}, \quad \operatorname{dim}\left((J L-i \mu) E_{i \mu}\right) \leq 2\left(k^{\leq 0}(i \mu)-n^{-}\left(\left.L\right|_{E^{1}}\right)\right) \\
\operatorname{dim} E^{G} \leq 3\left(k^{\leq 0}(i \mu)-\operatorname{dim} E^{D}-n^{-}\left(\left.L\right|_{E^{1}}\right)\right)
\end{gathered}
$$

and $L$ and $J L$ take the block forms on $E$

$$
\langle L \cdot, \cdot\rangle \longleftrightarrow\left(\begin{array}{ccc}
0 & 0 & 0 \\
0 & L_{1} & 0 \\
0 & 0 & L_{G}
\end{array}\right), \quad J L \longleftrightarrow\left(\begin{array}{ccc}
A_{D} & A_{D 1} & A_{D G} \\
0 & i \mu & 0 \\
0 & 0 & A_{G}
\end{array}\right)
$$

where all blocks are bounded operators and $L_{1}$ and $L_{G}$ are non-degenerate. Moreover, $\operatorname{ker}\left(A_{G}-i \mu\right) \subset\left(A_{G}-i \mu\right) E_{G}$.

Proof. Again, to apply previous results directly it would be easier to consider the complexifications of real Hilbert spaces

$$
I_{i \mu} \triangleq E_{i \mu}+E_{-i \mu}=\left\{u+\bar{v} \mid u, v \in E_{i \mu}\right\}
$$

Due to Lemma 3.5, $\left.J L\right|_{I_{i \mu}}$ is bounded with

$$
L\left(I_{i \mu}\right) \subset D(J), J L\left(I_{i \mu}\right) \subset I_{i \mu}, \sigma\left(\left.J L\right|_{I_{i \mu}}\right)=\{ \pm i \mu\}
$$

We split the spaces by starting with

$$
\begin{aligned}
& E^{D}=\left\{u \in E_{i \mu} \mid\langle L u, v\rangle=0, \forall v \in E_{i \mu}\right\}, I^{D}=\left\{u+\bar{v} \mid u, v \in E^{D}\right\} \\
& E_{i \mu}^{N D}=\left\{u \in E_{i \mu} \mid(u, v)=0, \forall v \in E^{D}\right\}, I^{N D}=\left\{u+\bar{v} \mid u, v \in E_{i \mu}^{N D}\right\}
\end{aligned}
$$

From the anti-symmetry of $J L$ with respect to $\langle L \cdot, \cdot\rangle$ and the invariance of $E_{i \mu}$ along with (12.12), we have $J L\left(I^{D}\right) \subset I^{D}$. In the splitting $I_{i \mu}=I^{D} \oplus I^{N D},\langle L \cdot, \cdot\rangle$ and $J L$ can be represented in the following block forms

$$
\langle L \cdot, \cdot\rangle \longleftrightarrow\left(\begin{array}{cc}
0 & 0 \\
0 & L_{N D}
\end{array}\right), \quad J L \longleftrightarrow\left(\begin{array}{cc}
A_{D} & A_{D, N D} \\
0 & A_{N D}
\end{array}\right)
$$

where all blocks are bounded real (satisfying (12.12)) operators. In particular, $\operatorname{ker}\left(\left.L\right|_{I_{i \mu}}\right)=I^{D}$ and $I_{i \mu}=I^{D} \oplus I^{N D}$ and thus Lemma 12.2 implies that $L_{N D}$ : $I^{N D} \rightarrow\left(I^{N D}\right)^{*}$ is an isomorphism. The anti-symmetry of $J L$ with respect to $\langle L \cdot, \cdot\rangle$ yields $L_{N D} A_{N D}+A_{N D}^{*} L_{N D}=0$. Therefore,

$$
J_{N D}=A_{N D} L_{N D}^{-1}:\left(I^{N D}\right)^{*} \rightarrow I_{N D}
$$

is an anti-symmetric bounded operator satisfying $A_{N D}=J_{N D} L_{N D}$. Clearly, the combination $\left(I^{N D}, L_{N D}, J_{N D}\right)$ satisfies (H1-3) with the non-degenerate $L_{N D}$. Moreover, $\sigma\left(J_{N D} L_{N D}\right)=\{ \pm i \mu\}$ with the eigenspace of $i \mu$ given by $E_{i \mu}^{N D}$ where $\left\langle L_{N D^{\cdot}}, \cdot\right\rangle$ is also non-degenerate. Therefore, we may apply Lemma 8.3 (with $X$ and $E$ replaced by $I^{N D}$ and $E_{i \mu}^{N D}$, respectively) to obtain the splitting $E_{i \mu}^{N D}=E^{1} \oplus E^{G}$ and the desired block forms of $L$ and $J L$ follow. The desired estimate on $\operatorname{dim} E^{G}$ is obtained by noting

$$
\begin{equation*}
k^{\leq 0}(i \mu)=n^{-}\left(\left.L\right|_{E_{i \mu}^{N D}}\right)+\operatorname{dim} E^{D} . \tag{8.5}
\end{equation*}
$$

Moreover, according to Lemma 8.3, we have

$$
\operatorname{dim}\left(A_{N D}-i \mu\right) E_{i \mu}^{N D} \leq 2\left(n^{-}\left(\left.L\right|_{E_{i \mu}^{N D}}\right)-n^{-}\left(\left.L\right|_{E^{1}}\right)\right)
$$

Along with (8.5) and the block form of $J L$, it implies

$$
\begin{aligned}
& \operatorname{dim}(J L-i \mu) E_{i \mu} \leq \operatorname{dim} E^{D}+\operatorname{dim}\left(A_{N D}-i \mu\right) E_{i \mu}^{N D} \\
\leq & \operatorname{dim} E^{D}+2\left(n^{-}\left(\left.L\right|_{E_{i \mu}^{N D}}\right)-n^{-}\left(\left.L\right|_{E^{1}}\right)\right) \leq 2\left(k^{\leq 0}(i \mu)-n^{-}\left(\left.L\right|_{E^{1}}\right)\right)
\end{aligned}
$$

which finishes the proof.
Proof of Proposition 2.1 and Proposition 2.2. What remains to be proved in these two propositions can be obtained in a similar framework and we complete their proofs together here.

Let $X_{ \pm}$be given by Lemma 12.4 and $X_{1}=X_{-} \oplus X_{+}$. Clearly, $X=X_{0} \oplus X_{1}$, where $X_{0}=\operatorname{ker} L$, with the associated projections $P_{X_{0,1}}$. According to Lemma $12.3,\langle L \cdot, \cdot\rangle$ and $J L$ take the following block forms

$$
\langle L \cdot, \cdot\rangle \longleftrightarrow\left(\begin{array}{cc}
0 & 0 \\
0 & L_{X_{1}}
\end{array}\right), \quad J L \longleftrightarrow\left(\begin{array}{cc}
0 & A_{1} \\
0 & J_{X_{1}} L_{X_{1}}
\end{array}\right)
$$

where $A_{1}: X_{1} \rightarrow \operatorname{ker} L$ is bounded and $L_{X_{1}}=i_{X_{1}}^{*} L i_{X_{1}}: X_{1} \rightarrow X_{1}^{*}$ and $J_{X_{1}}=$ $P_{X_{1}} J P_{X_{1}}^{*}$. Moreover, Lemmas 12.3 and 12.4 imply that $\left(X_{1}, L_{X_{1}}, J_{X_{1}}\right)$ satisfies assumptions (H1-3) with the isomorphic $L_{X_{1}}$ and $n^{-}\left(L_{X_{1}}\right)=n^{-}(L)$. For any eigenvalue $i \mu \in i \mathbf{R}$, let $E_{i \mu}^{1}$ be the subspace of generalized eigenvectors of $i \mu$ for $J_{X_{1}} L_{X_{1}}$, possibly $\{0\}$ if $\mu=0$. From Lemma 3.5 and 8.4 , for some $K>0$,

$$
E_{i \mu}^{1}=\operatorname{ker}\left(J_{X_{1}} L_{X_{1}}-i \mu\right)^{K}, \quad \operatorname{dim}\left(J_{X_{1}} L_{X_{1}}-i \mu\right) E_{i \mu}^{1} \leq 2 n^{\leq 0}\left(\left.L_{X_{1}}\right|_{E_{i \mu}^{1}}\right)
$$

For any integer $k>0,(J L-i \mu)^{k}$ takes the block form

$$
(J L-i \mu)^{k} \longleftrightarrow\left(\begin{array}{cc}
(-i \mu)^{k} & A_{k} \\
0 & \left(J_{X_{1}} L_{X_{1}}-i \mu\right)^{k}
\end{array}\right)
$$

where the linear operator $A_{k}: X_{1} \rightarrow \operatorname{ker} L$ can be computed inductively

$$
A_{k+1}=(-i \mu)^{k} A_{1}+A_{k}\left(J_{X_{1}} L_{X_{1}}-i \mu\right), \quad D\left(\left(J_{X_{1}} L_{X_{1}}-i \mu\right)^{k}\right) \subset D\left(A_{k+1}\right)
$$

It is straightforward to show

$$
\begin{equation*}
u \in E_{i \mu} \Longleftrightarrow P_{X_{1}} u \in E_{i \mu}^{1} \text { and }(-i \mu)^{K} P_{X_{0}} u+A_{K} P_{X_{1}} u=0 \tag{8.6}
\end{equation*}
$$

We first consider $\mu \neq 0$. We obtain from (8.6)

$$
E_{i \mu}=\left\{u-(-i \mu)^{-K} A_{K} u \mid u \in E_{i \mu}^{1}\right\}
$$

i.e. vectors in $E_{i \mu}$ are determined only by their $X_{1}$-component. From Lemma 3.5 and Remark 3.1, $E_{i \mu}^{1}$ and $E_{i \mu}$ are both subspaces. Therefore, $A_{K}$ is a bounded operator. Since

$$
\left\langle L\left(u-(-i \mu)^{-K} A_{K} u\right), v-(-i \mu)^{-K} A_{K} v\right\rangle=\langle L u, v\rangle, \forall u, v \in E_{i \mu}^{1}
$$

we obtain from Lemma 8.4

$$
\begin{aligned}
\operatorname{dim}(J L-i \mu) E_{i \mu} & =\operatorname{dim}\left(J_{X_{1}} L_{X_{1}}-i \mu\right) E_{i \mu}^{1} \\
& \leq 2 n^{\leq 0}\left(\left.L_{X_{1}}\right|_{E_{i \mu}^{1}}\right)=2 n^{\leq 0}\left(L_{E_{i \mu}}\right)=2 k^{\leq 0}(i \mu)
\end{aligned}
$$

This proves the desired estimate on $\operatorname{dim}(J L-i \mu) E_{i \mu}$ in Proposition 2.1. Along with Lemma 3.5, it completes the proof of Proposition 2.1 in the case of $\mu \neq 0$.

To prove Proposition 2.2, let $E_{i \mu}^{1}=\tilde{E}^{D} \oplus \tilde{E}^{1} \oplus \tilde{E}^{G}$ where these subspaces are given by Lemma 8.4 for $J_{X_{1}} L_{X_{1}}$. Let

$$
E^{D, 1, G}=\left\{u-(-i \mu)^{-K} A_{K} u \mid u \in \tilde{E}^{D, 1, G}\right\}
$$

It is easy to verify that they satisfy the properties in Proposition 2.2. Since $\operatorname{dim} E^{G}<\infty$, the 'good' basis of $E^{G}$ has been constructed in the finite dimensional cases in Chapter 4 and the proof of Proposition 2.2 is complete.

For $\mu=0$, it is easy to see from the above block forms

$$
E_{0}=X_{0} \oplus E_{0}^{1}
$$

Therefore, we have

$$
(J L)^{2} E_{0}=(J L)^{2} E_{0}^{1}=J L\left(P_{X_{1}} J L E_{0}^{1}\right)=J L\left(J_{X_{1}} L_{X_{1}} E_{0}^{1}\right)
$$

which along with Lemma 8.4 implies

$$
\operatorname{dim}(J L)^{2} E_{0} \leq \operatorname{dim} J_{X_{1}} L_{X_{1}} E_{0}^{1} \leq 2 k_{0}^{\leq 0}
$$

This completes the proof of Proposition 2.1 in the case of $\mu=0$.
To prove Proposition 2.2, let $E_{0}^{1}=E^{D} \oplus E^{1} \oplus E^{G}$ where these subspaces are given by Lemma 8.4 for $J_{X_{1}} L_{X_{1}}$ and $\mu=0$. It is easy to verify that they satisfy the properties in Proposition 2.2. Again since $\operatorname{dim} E^{G}<\infty$, the 'good' basis of $E^{G}$ has been constructed in the finite dimensional cases in Chapter 4 and the proof of Proposition 2.2 is complete.

REMARK 8.2. In the case of $\mu=0$, we can not replace $(J L)^{2} E_{0}$ by $J L E_{0}$, as seen from the following counterexample. Consider $X=Y \oplus Y \oplus \mathbf{R}^{\mathbf{2}}$ where $Y$ is any Hilbert space. Let

$$
J=\left(\begin{array}{cccc}
0 & I & 0 & 0 \\
-I & 0 & 0 & 0 \\
0 & 0 & 0 & -1 \\
0 & 0 & 1 & 0
\end{array}\right), \quad L=\left(\begin{array}{cccc}
0 & 0 & 0 & 0 \\
0 & I & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & -1
\end{array}\right), \quad J L=\left(\begin{array}{cccc}
0 & I & 0 & 0 \\
0 & 0 & 0 & 0 \\
0 & 0 & 0 & 1 \\
0 & 0 & 1 & 0
\end{array}\right)
$$

It is clear that $k_{0}^{<0}=0$, $\operatorname{ker} L=X_{0}=Y \oplus\{0\} \oplus\left\{(0,0)^{T}\right\}, E_{0}=Y \oplus Y \oplus\left\{(0,0)^{T}\right\}$, and $\operatorname{dim} J L E_{0}=\operatorname{dim} \operatorname{ker} L=\operatorname{dim} Y$.

### 8.3. Subspace of generalized eigenvectors $E_{0}$ and index $k_{0}^{\leq 0}$

In this Section we prove Propositions 2.7, 2.8, Lemma 2.1 and Corollary 2.3, 2.4 on the subspace $E_{0}$ and the non-positive index $k_{0}^{\leq 0}$ for the eigenvalue 0 .

Proof of Proposition 2.7. According to Corollary 12.1, $L J: D(J) \rightarrow X$ is closed and thus $L J L$ is also closed. Therefore, $(J L)^{-1}(\operatorname{ker} L)=\operatorname{ker}(L J L)$ is also closed.

Since $(J L)^{-1}(\operatorname{ker} L) \subset E_{0}$, due to the hyperbolicity of $J L$ on $X_{5,6}$, we have $(J L)^{-1}(\operatorname{ker} L) \subset \oplus_{j=0}^{4} X_{j}$, where the decomposition of $X=\operatorname{ker} L \oplus \oplus_{j=1}^{6} X_{j}$ is given in Theorem 2.1. Let

$$
S=(J L)^{-1}(\operatorname{ker} L) \cap \oplus_{j=1}^{4} X_{j}
$$

Since $X_{0}=\operatorname{ker} L \subset(J L)^{-1}(\operatorname{ker} L)$, we have $\operatorname{ker} L \oplus S=(J L)^{-1}(\operatorname{ker} L) \subset E_{0}$. Therefore, from the definition of $k_{0}^{\leq 0}$ it is clear $k_{0}^{\leq 0} \geq n_{0}=n^{\leq 0}\left(\left.L\right|_{S}\right)$ and we only need to prove (ii) of Proposition 2.7.

Assume in addition that $\langle L \cdot, \cdot\rangle$ is non-degenerate on $(J L)^{-1}(\operatorname{ker} L) / \operatorname{ker} L$. We claim

$$
\begin{equation*}
E_{0}=(J L)^{-1}(\operatorname{ker} L) \tag{8.7}
\end{equation*}
$$

In fact, suppose $u \in E_{0} \backslash\left((J L)^{-1}(\operatorname{ker} L)\right)$. There exists $m>0$ such that

$$
\begin{aligned}
& u_{1} \triangleq(J L)^{m-1} u \notin(J L)^{-1}(\operatorname{ker} L) \\
& u_{0} \triangleq J L u_{1}=(J L)^{m} u \in(J L)^{-1}(\operatorname{ker} L) \backslash \operatorname{ker} L
\end{aligned}
$$

It follows that, for any $v \in(J L)^{-1}(\operatorname{ker} L)$,

$$
J L v \in \operatorname{ker} L \Longrightarrow\left\langle L u_{0}, v\right\rangle=\left\langle L(J L) u_{1}, v\right\rangle=-\left\langle L u_{1}, J L v\right\rangle=0
$$

The existence of such $u_{0}$ would imply $\langle L \cdot, \cdot\rangle$ is degenerate on $(J L)^{-1}(\operatorname{ker} L) / \operatorname{ker} L$, contradictory to our assumption. Therefore, (8.7) is proved and consequently we obtain from the definition of $k_{0}^{\leq 0}$ that

$$
k_{0}^{\leq 0}=n^{\leq 0}\left(\left.\langle L \cdot, \cdot\rangle\right|_{(J L)^{-1}(\operatorname{ker} L) / \operatorname{ker} L}\right)=n^{-}\left(\left.\langle L \cdot, \cdot\rangle\right|_{(J L)^{-1}(\operatorname{ker} L) / \operatorname{ker} L}\right)
$$

due to the non-degeneracy assumption. This completes the proof of the proposition.
We will prove Lemma 2.1, Proposition 2.8, and Corollary 2.3 and 2.4 in the rest of the section. We first observe that it is straightforward to show $\langle L u, v\rangle=0$, for any $u \in \operatorname{ker}(J L)$ and $v \in R(J)$. Through a density argument, we obtain

$$
\begin{equation*}
\langle L u, v\rangle=0, \quad \forall u \in \operatorname{ker}(J L), v \in \overline{R(J)} \tag{8.8}
\end{equation*}
$$

Throughout the rest of this section, let $S_{1}, S_{2}, S^{\#}$ be defined as in Corollaries 2.3 and 2.4, i.e.

$$
\operatorname{ker}(J L)=\operatorname{ker} L \oplus S_{1}, \quad \overline{R(J)}=(\overline{R(J)} \cap \operatorname{ker} L) \oplus S^{\#}
$$

and

$$
\overline{R(J)} \cap(J L)^{-1}(\operatorname{ker} L)=S_{2} \oplus(\overline{R(J)} \cap \operatorname{ker} L)
$$

Lemma 8.5. Suppose $\langle L \cdot, \cdot\rangle$ is non-degenerate on $S^{\#}$, then it is also nondegenerate on $S_{1}$ and moreover,

$$
\begin{equation*}
X=\operatorname{ker}(J L) \oplus S^{\#}=\operatorname{ker} L \oplus S_{1} \oplus S^{\#} \tag{8.9}
\end{equation*}
$$

Proof. The non-degeneracy of $\langle L \cdot, \cdot\rangle$ on $S^{\#}$ implies the non-degeneracy of $L_{S \#}: S^{\#} \rightarrow\left(S^{\#}\right)^{*}$, which is defined in (12.1). For any $u \in X$, as in the proof of Lemma 12.2, let

$$
u^{\#}=L_{S \#}^{-1} i_{S \#}^{*} L u \in S^{\#}
$$

which satisfies

$$
\left\langle L u_{1}, v\right\rangle=0, \forall v \in S^{\#}, \quad \text { where } u_{1}=u-u^{\#}
$$

By the definition of $S^{\#}$, we also have

$$
\left\langle L u_{1}, v\right\rangle=0, \forall v \in \overline{R(J)}
$$

Since $J^{*}=-J$, we obtain

$$
L u_{1} \in \operatorname{ker} J^{*}=\operatorname{ker} J \Longrightarrow u_{1} \in \operatorname{ker}(J L)=\operatorname{ker} L \oplus S_{1}
$$

Therefore, $u=u_{1}+u^{\#} \in \operatorname{ker}(J L)+S^{\#}$ and thus $X=\operatorname{ker}(J L)+S^{\#}$.
For any $u \in S^{\#} \cap \operatorname{ker}(J L)$, from (8.8) we obtain $\langle L u, v\rangle=0$, for any $v \in$ $\operatorname{ker}(J L)+\overline{R(J)} \supset \operatorname{ker}(J L)+S^{\#}=X$. Therefore, $u \in \operatorname{ker} L$. Since $u \in S^{\#} \cap \operatorname{ker} L=$ $\{0\}$, we have $u=0$ and thus $X=\operatorname{ker}(J L) \oplus S^{\#}=\operatorname{ker} L \oplus S_{1} \oplus S^{\#}$.

From Lemma 12.2, $\langle L \cdot, \cdot\rangle$ is non-degenerate on $S^{\#} \oplus S_{1}$. Since it is also assumed to be non-degenerate on $S^{\#}$, the non-degeneracy of $\langle L \cdot, \cdot\rangle$ on $S_{1}$ follows from the $L$-orthogonality (8.8) between $S_{1}$ and $S^{\#}$.

Lemma 8.6. Suppose $\langle L \cdot, \cdot\rangle$ is non-degenerate on $S_{1}$, then it is also nondegenerate on $S^{\#}$.

Proof. Like in the proof of the previous lemma, the non-degeneracy of $\langle L \cdot, \cdot\rangle$ on $S_{1}$ implies the non-degeneracy of $L_{S_{1}}: S_{1} \rightarrow S_{1}^{*}$. For any $u \in X$, as in the proof of Lemma 12.2, let

$$
u_{1}=L_{S_{1}}^{-1} i_{S_{1}}^{*} L u \in S_{1}
$$

which satisfies

$$
\left\langle L u_{*}, v\right\rangle=0, \forall v \in \operatorname{ker}(J L)=\operatorname{ker} L \oplus S_{1}, \quad \text { where } u_{*}=u-u_{1}
$$

Since $J L=-(L J)^{*}$, we obtain $L u_{*} \in \overline{R(L J)}$.
Claim: $\overline{R(L J)}=L\left(S^{\#}\right)$. In fact, it is easy to see $L\left(S^{\#}\right) \subset L(\overline{R(J)}) \subset \overline{R(L J)}$ due to the boundedness of $L$. In the following we will prove that $\overline{R(L J)} \subset L\left(S^{\#}\right)$. Let $y \in \overline{R(L J)}$, there exists a sequence $y_{n}=L J x_{n}$ such that $y_{n} \rightarrow y$ as $n \rightarrow+\infty$. Since $\overline{R(J)}=\operatorname{ker} L \oplus S^{\#}$, let $J x_{n}=z_{n, 0}+z_{0, \#}$ where $z_{n, 0} \in \operatorname{ker} L$ and $z_{n, \#} \in S^{\#}$. As $y_{n}=L J x_{n}=L z_{n, \#} \rightarrow y$ and the non-degeneracy assumption of $\langle L \cdot, \cdot\rangle$ on $S^{\#}$ implies that $\left.L\right|_{S^{\#}}: S^{\#} \rightarrow L\left(S^{\#}\right)$ is an isomorphism, we obtain that $\left\{z_{n, \#}\right\}$ is a Cauchy sequence. Let $z_{n, \#} \rightarrow z_{\#} \in S^{\#}$ and then $y=L z_{\#} \in L\left(S^{\#}\right)$. The claim is proved.

We can now finish the proof of the lemma. Since we have proved

$$
L\left(u-u_{1}\right)=L u_{*} \in \overline{R(L J)}=L\left(S^{\#}\right)
$$

there exists $u_{\#} \in S^{\#}$ such that $L\left(u-u_{1}\right)=L u_{\#}$. Let $u_{0}=u-u_{1}-u_{\#}$. Clearly, $u_{0} \in$ $\operatorname{ker} L$. Therefore, $u=u_{0}+u_{1}+u_{\#}$ and thus $X=\operatorname{ker} L \oplus S_{1} \oplus S^{\#}=\operatorname{ker}(J L) \oplus S^{\#}$. The proof of $\operatorname{ker}(J L) \cap S^{\#}=\{0\}$ and consequently the non-degeneracy of $\langle L \cdot, \cdot\rangle$ on $S_{1}$ is the same as in the proof of the last lemma.

The conclusion in Lemma 2.1 is already contained in the above lemmas.
Proof of Proposition 2.8 and equivalently Corollary 2.4. The property $X=\overline{\operatorname{ker}(J L)}+\overline{R(J)}$ is a direct consequence of (8.9). Along with (8.8), it also implies $\overline{R(J)} \cap \operatorname{ker}(J L)=\overline{R(J)} \cap$ ker $L \subset$ ker $L$.

From the $L$-orthogonality (8.8), the decomposition (8.9), and the non-degeneracy of $\langle L \cdot, \cdot\rangle$ on $S_{1}, S^{\#}$, and $S_{1} \oplus S^{\#}$, we immediately obtain $n^{-}=n^{-}\left(\left.L\right|_{S_{1}}\right)+n^{-}\left(\left.L\right|_{S^{\#}}\right)$.

From the decomposition (8.9) and the definitions of $S_{1}$ and $S_{2}$, we have

$$
(J L)^{-1} \operatorname{ker} L=\operatorname{ker} L \oplus S_{1} \oplus S_{2}
$$

Therefore, $k_{0}^{\leq 0} \geq n^{-}\left(\left.L\right|_{S_{1}}\right)+n^{\leq 0}\left(\left.L\right|_{S_{2}}\right)$ follows from Proposition 2.7.
Finally, let us assume, in addition, that $\langle L \cdot, \cdot\rangle$ is non-degenerate on $S_{2}$. Immediately we have the non-degeneracy of $\langle L \cdot, \cdot\rangle$ on $(J L)^{-1}(\operatorname{ker} L) / \operatorname{ker} L$ and Proposition 2.7 implies $k_{0}^{\leq 0}=n^{-}\left(\left.L\right|_{S_{1}}\right)+n^{\leq 0}\left(\left.L\right|_{S_{2}}\right)$. The proof is complete.

### 8.4. Non-degeneracy of $\langle L \cdot, \cdot\rangle$ on $E_{i \mu}$

In Proposition 2.2, the presence of the subspace $E^{D} \subset E_{i \mu}$ is due to the possible degeneracy of $\langle L \cdot, \cdot\rangle$ on $E_{i \mu}$. Otherwise the statement of the proposition would be much more clean and some results can be improved. However, in case when $i \mu$ is not isolated in $\sigma(J L)$, it is indeed possible that $\langle L \cdot, \cdot\rangle$ degenerates on $E_{i \mu}$ even if it is non-degenerate on $X$.

Example of degenerate $\langle L \cdot, \cdot\rangle$ on $E_{i \mu}$. Consider $X=\mathbf{R}^{2 n} \oplus \mathbf{R}^{2 n} \oplus X_{1}$, where $X_{1}$ is a Hilbert space. Here we identify Hilbert spaces and their dual spaces via Riesz Representation Theorem. Let $\mu \in \mathbf{R}$ and

- $A: X_{1} \supset D(A) \rightarrow X_{1}$ be an anti-self-adjoint operator such that $i \mu \in \sigma(A)$ is not an eigenvalue;
- $A_{1}: \mathbf{R}^{2 n} \rightarrow X_{1}$ such that $\operatorname{ker} A_{1}=\{0\}$ and, after the complexification of $A$ and $A_{1}$ into complex linear operators, $R\left(A_{1}\right) \cap R(A \pm i \mu)=\{0\}$, which is possible due to the spectral assumption on $A$; and
- $J=\left(\begin{array}{ccc}0 & J_{2 n} & 0 \\ J_{2 n} & J_{2 n} & -B^{-1} A_{1}^{*} \\ 0 & A_{1} B^{-1} & A\end{array}\right), L=\left(\begin{array}{ccc}0 & B & 0 \\ B & 0 & 0 \\ 0 & 0 & I_{X_{1}}\end{array}\right)$, where $B_{2 n \times 2 n}$ is any symmetric matrix and $J_{2 n}=\left(\begin{array}{cc}0 & -I_{n \times n} \\ I_{n \times n} & 0\end{array}\right)$.
One may compute

$$
J L=\left(\begin{array}{ccc}
J_{2 n} B & 0 & 0 \\
J_{2 n} B & J_{2 n} B & -B^{-1} A_{1}^{*} \\
A_{1} & 0 & A
\end{array}\right)
$$

Lemma 8.7. For any integer $k>0$,

$$
\operatorname{ker}(J L-i \mu)^{k}=\left\{(0, x, 0)^{T} \mid x \in \operatorname{ker}\left(J_{2 n} B-i \mu\right)^{k}\right\} \subset \mathbf{R}^{2 n} \times \mathbf{R}^{2 n} \times X_{1}
$$

Consequently, $\langle L \cdot, \cdot\rangle$ vanishes on $E_{ \pm i \mu}$.
REmARK 8.3. The embedding from $\mathbf{R}^{2 n}$ to $\{0\} \times \mathbf{R}^{2 n} \times\{0\} \subset X$ - an invariant subspace under $J L$, serves as a similarity transformation between the $2 n$-dim Hamiltonian operator $J_{2 n} B$ and the restriction of the infinite dimensional one JL. If $i \mu \in \sigma\left(J_{2 n} B\right)$, then $J_{2 n} B$ and $J L$ have exactly the same structures on the subspaces $E_{i \mu}\left(J_{2 n} B\right)$ and $E_{i \mu}$ of generalized eigenvectors of $i \mu$. However, the energy structure is completely destroyed. Namely the $2 n$-dim Hamiltonian operator $J_{2 n} B$ has a non-trivial energy $\langle B \cdot, \cdot\rangle$ while the energy $\langle L \cdot, \cdot\rangle$ of $J L$ vanishes completely on $\mathbf{R}^{2 n}$ to $\{0\} \times \mathbf{R}^{2 n} \times\{0\} \subset X$.

Proof. Using the invariance under $J L$ of $\{0\} \times \mathbf{R}^{2 n} \times\{0\}$ and $\{0\} \times \mathbf{R}^{2 n} \times X_{1}$, it is easy to compute inductively

$$
(J L-i \mu)^{k}=\left(\begin{array}{ccc}
\left(J_{2 n} B-i \mu\right)^{k} & 0 & 0 \\
A_{21} & \left(J_{2 n} B-i \mu\right)^{k} & A_{23} \\
A_{31} & 0 & (A-i \mu)^{k}
\end{array}\right)
$$

where

$$
A_{31}=\Sigma_{l=0}^{k-1}(A-i \mu)^{l} A_{1}\left(J_{2 n} B-i \mu\right)^{k-1-l}
$$

Let $P_{1,2,3}$ denote the projections from $X$ to its components. For any $u=\left(x_{1}, x_{2}, v\right)^{T} \in$ $X$, we have

$$
\begin{aligned}
P_{3}(J L-i \mu)^{k} u= & A_{31} x_{1}+(A-i \mu)^{k} v \\
= & A_{1}\left(J_{2 n} B-i \mu\right)^{k-1} x_{1}+(A-i \mu)\left((A-i \mu)^{k-1} v\right. \\
& \left.\quad+\Sigma_{l=1}^{k-1}(A-i \mu)^{l-1} A_{1}\left(J_{2 n} B-i \mu\right)^{k-1-l} x_{1}\right)
\end{aligned}
$$

Suppose $P_{3}(J L-i \mu)^{k} u=0$. Since $A_{1}$ and $A-i \mu$ are both one-to-one and $R\left(A_{1}\right) \cap R(A-i \mu)=\{0\}$, we obtain

$$
\left(J_{2 n} B-i \mu\right)^{k-1} x_{1}=0,(A-i \mu)^{k-1} v+\Sigma_{l=1}^{k-1}(A-i \mu)^{l-1} A_{1}\left(J_{2 n} B-i \mu\right)^{k-1-l} x_{1}=0 .
$$

Let $m \in[0, k-1]$ be the minimal non-negative integer satisfying $\left(J_{2 n} B-i \mu\right)^{m} x_{1}=0$. If $m \geq 1$, from the definition of $m$, the above second equality and the injectivity of $(A-i \mu)^{k-1-m}$ imply

$$
\begin{aligned}
0= & (A-i \mu)^{m} u+\Sigma_{l=k-m}^{k-1}(A-i \mu)^{l+m-k} A_{1}\left(J_{2 n} B-i \mu\right)^{k-1-l} x_{1} \\
= & (A-i \mu)\left((A-i \mu)^{m-1} v+\Sigma_{l=k-m+1}^{k-1}(A-i \mu)^{l+m-k-1} A_{1}\left(J_{2 n} B-i \mu\right)^{k-1-l} x_{1}\right) \\
& +A_{1}\left(J_{2 n} B-i \mu\right)^{m-1} x_{1}
\end{aligned}
$$

Again since $A_{1}$ and $A-i \mu$ are both one-to-one and $R\left(A_{1}\right) \cap R(A-i \mu)=\{0\}$, we derive $\left(J_{2 n} B-i \mu\right)^{m-1} x_{1}=0$ which contradicts the definition of $m$. Therefore, $m=0$, that is, $x_{1}=0$. Due to the injectivity of $(A-i \mu)^{k}$, it implies $v=0$ as well.

Suppose $u \in \operatorname{ker}(J L-i \mu)^{k}$, the above arguments imply $u=(0, x, 0)^{T}$ and the lemma follows immediately.

In the rest of this section we will prove that the degeneracy of $\langle L \cdot, \cdot\rangle$ may occur on $E_{i \mu}$ only if $i \mu \in \sigma(J L)$ is not an isolated spectral point.

Lemma 8.8. Assume (H1-3) and $i \mu \in \sigma(J L) \cap i \mathbf{R}$ is isolated in $\sigma(J L)$, then there exist closed subspaces $I^{i \mu}, E_{\&} \subset X$ such that
(i) $I^{i \mu}$ and $E_{\&}$ are complexifications of real subspaces of $X$, namely $u \in I^{i \mu}$ (or $E_{\&}$ ) if and only if $\bar{u} \in I^{i \mu}$ (or $E_{\&}$ ). Moreover, they are invariant under $J L$ and

$$
X=I^{i \mu} \oplus E_{\&}, \quad \sigma\left(\left.J L\right|_{I^{i \mu}}\right)=\{ \pm i \mu\}, \quad \sigma\left(\left.J L\right|_{E_{\&}}\right)=\sigma(J L) \backslash\{ \pm i \mu\}
$$

(ii) $\operatorname{ker} L \subset E_{\&}$ if $\mu \neq 0$ or $\operatorname{ker} L \subset I^{i \mu}$ if $\mu=0$.
(iii) $\langle L u, v\rangle=0$ for all $u \in I^{i \mu}$ and $v \in E_{\&}$. Moreover, $\langle L \cdot, \cdot\rangle$ is non-degenerate on quotient spaces $I^{i \mu} /\left(\operatorname{ker} L \cap I^{i \mu}\right)$ and $E_{\&} /\left(\operatorname{ker} L \cap E_{\&}\right)$.
Proof. Let $\Gamma \subset \mathbf{C} \backslash \sigma(J L)$ be a small circle, oriented counterclockwisely, enclosing $i \mu$ but no other elements in $\sigma(J L)$. Define the spectral projection and the eigenspaces

$$
P_{i \mu}=\frac{1}{2 \pi i} \oint_{\Gamma}(\lambda-J L)^{-1} d \lambda, \quad E^{i \mu}=P_{i \mu} X
$$

It is standard to verify that $P_{i \mu}$ is a bounded projection on $X$ satisfying

$$
J L P_{i \mu}=P_{i \mu} J L ; \sigma\left(\left.(J L)\right|_{E^{i \mu}}\right)=\{i \mu\} ; E^{i \mu} \subset D(J L) ; e^{t J L} E^{i \mu}=E^{i \mu}, \forall t \in \mathbf{R}
$$

By Lemma 3.6, $-i \mu \in \sigma(J L)$ is also an isolated point of $\sigma(J L)$. Let $P_{-i \mu}$ and $E^{-i \mu}$ be defined similarly. It is standard that $P_{i \mu} P_{-i \mu}=P_{-i \mu} P_{i \mu}=0$ and thus $P_{i \mu}+P_{-i \mu}$ is also a projection (or $P_{i \mu}$ instead if $\mu=0$ ). Define

$$
E_{\&}=\operatorname{ker}\left(P_{i \mu}+P_{-i \mu}\right), \quad I^{i \mu}=E^{i \mu}+E^{-i \mu}
$$

and we have

$$
\begin{equation*}
e^{t J L} E_{\&}=E_{\&}, \forall t \in \mathbf{R} ; \sigma\left(\left.(J L)\right|_{E_{\&}}\right)=\sigma(J L) \backslash\{ \pm i \mu\} ; X=I^{i \mu} \oplus E_{\&} \tag{8.10}
\end{equation*}
$$

Therefore, statements (i) and (ii) in the lemma follow from the standard spectral theory. The $L$-orthogonality between $I^{i \mu}$ and $E_{\&}$ follows from Lemma 6.2 where $\Omega$ can be taken as the union of the two small disks centered at $\pm i \mu$.

To complete the proof of the lemma, it suffices to prove the non-degeneracy of $\langle L \cdot, \cdot\rangle$ on $I^{i \mu} /\left(\operatorname{ker} L \cap I^{i \mu}\right)$ and $E_{\&} /\left(\operatorname{ker} L \cap E_{\&}\right)$. According to Lemma 12.3, $L_{I^{i \mu}}$ and $L_{E_{\&}}$ satisfy (H2). Therefore, either they are non-degenerate or have non-trivial kernels. Suppose there exists $u \in I^{i \mu}$ such that $\langle L u, v\rangle=0$ for all $v \in I^{i \mu}$. From
$X=I^{i \mu} \oplus E_{\&}$ and the $L$-orthogonality between $I^{i \mu}$ and $E_{\&}$, we obtain $\langle L u, v\rangle=0$ for all $v \in X$, which implies $u \in \operatorname{ker} L$. Therefore, $\langle L \cdot, \cdot\rangle$ is non-degenerate on $I^{i \mu} /\left(\operatorname{ker} L \cap I^{i \mu}\right)$. The proof of the non-degeneracy of $\langle L \cdot, \cdot\rangle$ on $E_{\&} /\left(\operatorname{ker} L \cap E_{\&}\right)$ is similar and thus we complete the proof of the lemma.

Notice that $I^{i \mu}$ is given not in terms of $E_{i \mu}$, but of $E^{i \mu}$ defined using spectral integrals. In the following we establish the relationship between $I^{i \mu}$ and the subspace $E_{i \mu}$ of generalized eigenvectors.

Lemma 8.9. It holds $I^{i \mu}=E_{i \mu}+E_{-i \mu}$.
Proof. Let $P^{\mu}: X \rightarrow I^{i \mu}$ be the projection associated to the $L$-orthogonal decomposition $X=I^{i \mu} \oplus E_{\&}$. Let $J^{\mu}=P^{\mu} J\left(P^{\mu}\right)^{*}$. As $I^{i \mu} \subset D(J L)$, Lemma 12.3 implies that $\left(I^{i \mu}, L_{I^{i \mu}}, J^{\mu}\right)$ satisfies assumptions (H1-3). The invariance of $I^{i \mu}$ under $J L$ implies $\left.J L\right|_{I^{i \mu}}=J^{\mu} L_{I^{i \mu}}$ and $\sigma\left(J^{\mu} L_{I^{i \mu}}\right)=\{ \pm i \mu\}$. Since $\pm i \mu \notin \sigma\left(\left.J L\right|_{E_{\&}}\right)$, we have $E_{ \pm i \mu} \subset I^{i \mu}$.

We apply Theorem 2.1 to $J L$ on $I^{i \mu}$, where there is no hyperbolic subspace, and obtain the decomposition of $I^{i \mu}$ into closed subspaces $I^{i \mu}=\Sigma_{j=0}^{4} X_{j}$, where $X_{0}=\operatorname{ker} L$ if $\mu=0$ or $X_{0}=\{0\}$ if $\mu \neq 0$. In this decomposition, $L_{I^{i \mu}}$ and $J L$ take the block forms

$$
J L \leftrightarrow\left(\begin{array}{ccccc}
0 & A_{01} & A_{02} & A_{03} & A_{04} \\
0 & A_{1} & A_{12} & A_{13} & A_{14} \\
0 & 0 & A_{2} & 0 & A_{24} \\
0 & 0 & 0 & A_{3} & A_{34} \\
0 & 0 & 0 & 0 & A_{4}
\end{array}\right), L_{I^{i \mu}} \leftrightarrow\left(\begin{array}{ccccc}
0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & B_{14} \\
0 & 0 & L_{X_{2}} & 0 & 0 \\
0 & 0 & 0 & L_{X_{3}} & 0 \\
0 & B_{14}^{*} & 0 & 0 & 0
\end{array}\right) .
$$

Note $L_{X_{3}} \geq \delta$ for some $\delta>0$ and $A_{2,3}$ are anti-self-adjoint with respect to the equivalent inner product $\mp\left\langle L_{X_{2,3}} \cdot, \cdot\right\rangle$ with $\sigma\left(A_{1,2,3,4}\right)=\{i \mu,-i \mu\}$.

In the case of $\mu=0$, the anti-self-adjoint operator $A_{2,3}$ must be $A_{2,3}=0$. Meanwhile all other finite dimensional diagonal blocks are also nilpotent. Therefore, it is straightforward to compute that $\left(\left.J L\right|_{I^{i \mu}}\right)^{k}=0$ for some integer $k>0$. Therefore, $I^{i \mu}$ consists of generalized eigenvectors only and $I^{i \mu}=E_{i \mu}$ in the case of $\mu=0$.

In the case of $\mu \neq 0, X_{0}=\{0\}$. Moreover, as $A_{3}$ is anti-self-adjoint with respect to the inner product $\left\langle L_{X_{3}} \cdot, \cdot\right\rangle$, we can further decompose $X_{3}$ into closed subspaces $X_{3}=X_{3+} \oplus X_{3-}$, where $X_{3 \pm}=\operatorname{ker}\left(A_{3} \pm i \mu\right)$, with associated projections $Q_{ \pm}: X_{3} \rightarrow X_{3 \pm}$. Accordingly $A_{3}=i \mu Q_{+}-i \mu Q_{-}$, which implies $A_{3}^{2}+\mu^{2}=0$. As $A_{1,2,4}$ are finite dimensional with the only eigenvalues $\pm i \mu$, we obtain that $\left(\left(\left.J L\right|_{I^{i \mu}}\right)^{2}+\mu^{2}\right)^{k}=0$ for some integer $k>0$. Rewrite it as

$$
(J L-i \mu)^{k}(J L+i \mu)^{k}=(J L+i \mu)^{k}(J L-i \mu)^{k}=0 \text { on } I^{i \mu} .
$$

Let $X_{ \pm}$be the invariant eigenspace of $\pm i \mu$ of $\left.J L\right|_{I^{i \mu}}$ defined via spectral integrals. We have $I^{i \mu}=X_{+} \oplus X_{-}$. As $J L \pm i \mu$ is an isomorphism from $X_{ \pm}$to itself, we obtain from the above identity that $X_{ \pm}=\operatorname{ker}(J L \mp i \mu)^{k}$. Therefore, $X_{ \pm}$are the subspaces of generalized eigenvectors of $\pm i \mu$ of $J L$, that is,

$$
I^{i \mu}=\operatorname{ker}(J L-i \mu)^{k} \oplus \operatorname{ker}(J L+i \mu)^{k}=E_{i \mu} \oplus E_{-i \mu}
$$

Finally, let $E_{\#}=E_{\&}$ if $\mu=0$ or $E_{\#}=E_{-i \mu} \oplus E_{\&}$ if $\mu \neq 0$. In the case of $\mu \neq 0$, Lemmas 3.1, 12.2, 12.3 and the non-degeneracy of $\langle L \cdot, \cdot\rangle$ on $I^{i \mu}$ imply that
$\langle L \cdot, \cdot\rangle$ is non-degenerate on $E_{i \mu}$ and $E_{\&} / \operatorname{ker} L$. This along with the above lemmas completes the proof of Proposition 2.3.

Based on Proposition 2.3, we are ready to prove Proposition 2.4.

## Proof of Proposition 2.4: Let

$$
\Lambda=\left\{0 \neq i \mu \in \sigma(J L) \cap i \mathbf{R} \mid k^{\leq 0}(i \mu)>0\right\}
$$

which is a finite set according to (2.13) of Theorem 2.3.
Let $i \mu \in \Lambda$. We have that $\langle L \cdot, \cdot\rangle$ is non-degenerate on $E_{i \mu}$ either by our assumption if $i \mu$ is not isolated in $\sigma(J L)$ or by Proposition 2.3 if $i \mu$ is isolated. From Proposition 2.2, we have the $L$-orthogonal and $J L$-invariant decomposition $E_{i \mu}=E_{i \mu}^{1} \oplus E_{i \mu}^{G}$, where $E_{i \mu}^{1} \subset \operatorname{ker}(J L-i \mu)$, $\operatorname{dim} E_{i \mu}^{G}<\infty$, and $\langle L \cdot, \cdot\rangle$ is non-degenerate on both $E_{i \mu}^{1}$ and $E_{i \mu}^{G}$. Let $E_{i \mu}^{1,-} \subset E_{i \mu}^{1}$ be a subspace such that $\operatorname{dim} E_{i \mu}^{1,-}=n^{-}\left(\left.L\right|_{E_{i \mu}^{-}}\right)$and $\langle L \cdot, \cdot\rangle$ is negative definite on $E_{i \mu}^{1,-}$. Let
$E_{i \mu}^{f i n i t e}=E_{i \mu}^{G} \oplus E_{i \mu}^{1,-}$, which satisfies $\operatorname{dim} E_{i \mu}^{\text {finite }}<\infty, n^{-}\left(\left.L\right|_{E_{i \mu}^{f i n i t e}}\right)=k^{\leq 0}(i \mu)$. Moreover, $J L\left(E_{i \mu}^{\text {finite }}\right)=E_{i \mu}^{\text {finite }}$ according to its construction.

If $0 \notin \sigma(J L)$, let $E_{0}^{\text {finite }}=\{0\}$ and we may skip to the next step to define $N$ and $M$. Otherwise, our assumption and Propositions 2.3, 2.2 imply an $L$-orthogonal decomposition $E_{0}=\operatorname{ker} L \oplus E_{0}^{1} \oplus E_{0}^{G}$, where

$$
J L\left(E_{0}^{1}\right) \subset \operatorname{ker} L, \operatorname{dim} E_{0}^{G}<\infty, J L\left(E_{0}^{G}\right) \subset E_{0}^{G} \oplus \operatorname{ker} L
$$

and $\langle L \cdot, \cdot\rangle$ is non-degenerate on both $E_{0}^{1}$ and $E_{0}^{G}$. Let $E_{0}^{1,-} \subset E_{0}^{1}$ be such that $\operatorname{dim} E_{0}^{1,-}=n^{-}\left(\left.L\right|_{E_{0}^{1}}\right)$ and $\langle L \cdot, \cdot\rangle$ is negative definite on $E_{0}^{1,-}$. Let $E_{0}^{\text {finite }}=E_{0}^{1,-} \oplus$ $E_{0}^{G}$, which satisfies

$$
\operatorname{dim} E_{0}^{\text {finite }}<\infty, n^{-}\left(\left.L\right|_{E_{0}^{\text {finite }}}\right)=k_{0}^{\leq 0}
$$

Let

$$
N=\left(\oplus_{\operatorname{Re} \lambda \neq 0} E_{\lambda}\right) \oplus\left(\oplus_{i \mu \in \Lambda} E_{i \mu}^{\text {finite }}\right) \oplus E_{0}^{\text {finite }}, \quad \tilde{M}=N^{\perp_{L}}=(N \oplus \operatorname{ker} L)^{\perp_{L}} .
$$

Clearly, $\operatorname{dim} N<\infty, n^{-}\left(\left.L\right|_{N}\right)=n^{-}(L)$ (due to (2.13) of Theorem 2.3), $\langle L \cdot, \cdot\rangle$ is non-degenerate on $N$, and $N \oplus \operatorname{ker} L$ is invariant under $J L$. Therefore, $\tilde{M}$ is also invariant under $J L, X=N \oplus \tilde{M}$ (due to Lemma 12.2), $\operatorname{ker} L \subset \tilde{M}$, and $n^{-}\left(\left.L\right|_{\tilde{M}}\right)=0$. Moreover, $N$ and $M$ are complexifications of real subspaces as $E_{\lambda}$ and $E_{\bar{\lambda}}$ have exactly the same structure. Let $M \subset \tilde{M}$ be any closed subspace such that $\tilde{M}=M \oplus \operatorname{ker} L$ and this completes the proof of proposition.

To end this chapter, we prove the decomposition result Proposition 2.5 for $L$ -self-adjoint operators.

Proof of Proposition 2.5: In order to apply the previous results, which have been given in the framework of real Hilbert spaces, to prove this proposition, we first convert it into a problem on real Hilbert spaces. Recall $(\cdot, \cdot)$ and $\langle\cdot, \cdot\rangle$ denote the complex inner product and the complex duality pair between $X^{*}$ and $X$, respectively. Let $X_{r}$ be the same set as $X$ but equipped with the real inner product $(u, v)_{r}=\operatorname{Re}(u, v)$. On $X_{r}$, the $i-$ multiplication $i: X \rightarrow X$ becomes a real linear isometry $i_{r}: X_{r} \rightarrow X_{r}$ with $i_{r}^{2}=-I$. Let $L_{r}: X_{r} \rightarrow X_{r}^{*}$ be the linear symmetry bounded operator defined as $\left\langle L_{r} u, v\right\rangle_{r}=\operatorname{Re}\langle L u, v\rangle$ where $\langle\cdot, \cdot\rangle_{r}$ denote
the real duality pair between $X_{r}^{*}$ and $X_{r}$. Subsequently, the non-degeneracy of $L$ yields the non-degeneracy of $L_{r}$. Accordingly, $A$ becomes a real linear operator $A_{r}: X_{r} \supset D\left(A_{r}\right) \rightarrow X_{r}$. The linearity of $L$ and $A$ implies that $i_{r} A_{r}=A_{r} i_{r}$ and $L_{r} i_{r}=-i_{r}^{*} L_{r}$. Finally, that $A$ is $L$-self-adjoint is translated to the $L_{r}$-selfadjointness of $A_{r}$, namely, $L_{r} A_{r}=A_{r}^{*} L_{r}$.

Define $J=i_{r} A_{r} L_{r}^{-1}: X_{r}^{*} \rightarrow X_{r}$. The $L_{r}$-self-adjointness of $A_{r}$ implies $J^{*}=$ $-J$ and thus $i_{r} A_{r}=J L_{r}$ with $\left(J, L_{r}, X_{r}\right)$ satisfying (H1-3). It is easy to prove

$$
\sigma\left(A_{r}\right)=\sigma(A) \subset \mathbf{R}, \quad \sigma\left(i_{r} A_{r}\right)=\left(i \sigma\left(A_{r}\right)\right) \cup\left(-i \sigma\left(A_{r}\right)\right)
$$

so the nonzero eigenvalues of $i_{r} A_{r}$ are isolated and $\operatorname{ker}\left(i_{r} A_{r}\right)=\operatorname{ker} A$. It is straightforward to deduce the non-degeneracy of $\left\langle L_{r} \cdot, \cdot\right\rangle_{r}$ on $\operatorname{ker}\left(i_{r} A_{r}\right)$ from our non-degeneracy assumption of $\langle L \cdot, \cdot\rangle$ on ker $A$, and thus (H4) is satisfied. Thus by Proposition 2.4, there exists a decomposition $X_{r}=\tilde{N} \oplus \tilde{M}$ such that $\tilde{N}$ and $\tilde{M}$ are $L_{r}$-orthogonal and invariant under $i_{r} A_{r}, \operatorname{dim} \tilde{N}<\infty$ and $\left.L_{r}\right|_{\tilde{M}}>0$, which also implies $L$ is uniformly positive on $\tilde{M}$. Let

$$
N=\tilde{N}+i_{r} \tilde{N}, \quad M=N^{\perp_{L_{r}}}=\left\{u \in X_{r} \mid\left\langle L_{r} u, v\right\rangle_{r}=0, \forall v \in N\right\} \subset \tilde{M}
$$

Clearly, $\operatorname{dim} N \leq 2 \operatorname{dim} \tilde{N}<\infty$ and $\left\langle L_{r} \cdot, \cdot\right\rangle$ is uniformly positive on $M$, thus so is $\langle L \cdot, \cdot\rangle$ on $M$. To complete the proof, we only need to show $N, M \subset X$ are $L$-orthogonal and invariant under $A$. We first consider the $L$-orthogonality which also involves the imaginary part of the quadratic form of $L$. Suppose there exist $u_{1}, u_{2}, \in \tilde{N}$ and $v \in M$ such that $\left\langle L\left(u_{1}+i u_{2}\right), v\right\rangle=R e^{i \theta} \neq 0$. It implies $\langle L u, v\rangle=$ $\left\langle L_{r} u, v\right\rangle_{r}=R \in \mathbf{R} \backslash\{0\}$, where

$$
u=(\cos \theta) u_{1}+(\sin \theta) u_{2}+i_{r}\left((\cos \theta) u_{2}+(\sin \theta) u_{1}\right) \in N
$$

which is a contradiction to the definitions of $N$ and $M$ and thus they are $L$ orthogonal. Secondly, $i_{r} A_{r}(\tilde{N}) \subset \tilde{N}, i_{r} A_{r}=A_{r} i_{r}$, and $i_{r}^{2}=-I$ imply $A_{r}\left(i_{r} \tilde{N}\right) \subset \tilde{N}$ and $A_{r} \tilde{N} \subset i_{r} \tilde{N}$. Therefore, $N$ is invariant under $A$. It along with the $L$-selfadjointness of $A$ also implies the invariance of $M$ and the proof of Proposition 2.5 is complete.

## CHAPTER 9

## Perturbations

In this chapter we study the robustness of the spectral properties of the Hamiltonian operator $J L$ under small perturbations preserving Hamiltonian structures. Consider

$$
u_{t}=J_{\#} L_{\#} u, \quad J_{\#}=J+J_{1}, \quad L_{\#}=L+L_{1}, \quad u \in X
$$

Unless otherwise specified, assumptions (A1-3) given in Section 2.5 are assumed throughout this chapter. We first prove

Lemma 9.1. Assumptions (A1-3) imply that there exists $\epsilon>0$ depending on $J$ and $L$ such that, if $\left|L_{1}\right| \leq \epsilon$, then $(\boldsymbol{H} 1-3)$ is satisfied by $J_{\#}$ and $L_{\#}$ and

$$
\operatorname{dim} \operatorname{ker} L_{\#} \leq \operatorname{dim} \operatorname{ker} L<\infty, \quad D\left(J_{\#} L_{\#}\right)=D(J L)
$$

Proof. It is obvious that $(\mathbf{H 1})$ is satisfied by $J_{\#}$. Let $X_{ \pm}$be the subspaces provided in (H2) satisfied by $L$. Clearly, we still have, for $\epsilon \ll 1$,

$$
\pm\left\langle L_{\#} u, u\right\rangle \geq \delta\|u\|^{2}, \quad \forall u \in X_{ \pm}
$$

for some $\delta>0$ independent of $\epsilon$. Let $X_{1}=X_{+} \oplus X_{-}$. Assumption (H2) for $L$ implies that $\langle L \cdot, \cdot\rangle$ restricted to $X_{1}$ is non-degenerate, i.e.

$$
L_{X_{1}}=i_{X_{1}}^{*} L i_{X_{1}}: X_{1} \rightarrow X_{1}^{*}
$$

defined as in (12.1), is an isomorphism. Therefore,

$$
L_{\#, X_{1}}=i_{X_{1}}^{*} L_{\#} i_{X_{1}}: X_{1} \rightarrow X_{1}^{*}
$$

as a small bounded perturbation of $L_{X_{1}}$, is also an isomorphism. Suppose $u=$ $u_{0}+u_{1} \in \operatorname{ker} L_{\#}$, where $u_{0} \in \operatorname{ker} L$ and $u_{1} \in X_{1}$, then we have

$$
0=L_{\#} u=L_{\#} u_{1}+L_{1} u_{0} \Longrightarrow u_{1}=-L_{\#, X_{1}}^{-1} i_{X_{1}}^{*} L_{1} u_{0}
$$

that is,
$\operatorname{ker} L_{\#} \subset Y, \quad$ where $Y=\operatorname{graph}(S)$ and $S=-L_{\#, X_{1}}^{-1} i_{X_{1}}^{*} L_{1}: \operatorname{ker} L \rightarrow X_{1}$.
Moreover, from the (12.2) type identity, it also holds that, for any $v \in \operatorname{ker} L$ and $u_{1} \in X_{1}$,

$$
\begin{aligned}
\left\langle L_{\#}(v+S v), u_{1}\right\rangle & =\left\langle L_{1} v, u_{1}\right\rangle-\left\langle L_{\#} L_{\#, X_{1}}^{-1} i_{X_{1}}^{*} L_{1} v, u_{1}\right\rangle \\
& =\left\langle L_{1} v, u_{1}\right\rangle-\left\langle i_{X_{1}}^{*} L_{1} v, u_{1}\right\rangle=0,
\end{aligned}
$$

that is, $Y$ and $X_{1}$ are $L_{\# \text {-orthogonal. }}$.
Since $\operatorname{dim} Y=\operatorname{dim} \operatorname{ker} L<\infty$ due to (A2), the quadratic form $\left\langle L_{\#} \cdot, \cdot\right\rangle$ restricted to $Y$ leads to a decomposition of $Y$

$$
Y=Y_{+} \oplus \operatorname{ker} L_{\#} \oplus Y_{-}
$$

where $\pm L_{\#}$ is positive on $Y_{ \pm}$. Let $X_{\# \pm}=X_{ \pm} \oplus Y_{ \pm}$, then

$$
X=X_{\#+} \oplus \operatorname{ker} L_{\#} \oplus X_{\#-}
$$

 are positive definite on $X_{\# \pm}$. Therefore, (H2) is satisfied.

Finally we prove (H3). Suppose $\gamma \in X^{*}$ and $\langle\gamma, u\rangle=0$ for all $u \in X_{\#+} \oplus$ $X_{\#-} \supset X_{+} \oplus X_{-}$. From (A1) which requires that (H1-3) being satisfied by $J$ and $L$, we have $\gamma \in D(J)=D\left(J_{\#}\right)$ as $J_{1}$ is assumed to be bounded.

Much as in Remark 2.2 by composing with the Riesz representation, we may treat $L_{\#}$ as a bounded symmetric operator on $X$ and then apply its spectral decomposition, a decomposition satisfying (H2) can be obtained much more easily. However, that decomposition may not satisfy (H3).

In Section 9.1, we will obtain the persistence of exponential trichotomy of the perturbed system. In Section 9.2, we will focus on purely imaginary spectral points of $\sigma(J L)$ and the possibility of bifurcation of unstable eigenvalues of $J_{\#} L_{\#}$. To start, following the standard procedure we show that assumption (A3) implies the convergence of the resolvents. Recall $\|\cdot\|_{G}$ denote the graph norm on $D(J L)$ and $|\cdot|_{G}$ the corresponding operator norm.

Lemma 9.2. Let $K \subset \mathbf{C} \backslash \sigma(J L)$ be compact, then there exist $C, \epsilon>0$ depending on $K$, $J$, and $L$, such that, for any $\lambda \in K$ and

$$
\left|J_{1}\right|,\left|J L_{1}\right|_{G} \leq \epsilon, \quad\left|L_{1}\right| \leq 1
$$

it holds that the densely defined closed operator $\lambda-J_{\#} L_{\#}: D(J L) \rightarrow X$ has a bounded inverse and

$$
\left|\left(\lambda-J_{\#} L_{\#}\right)^{-1}-(\lambda-J L)^{-1}\right| \leq C\left(\left|J_{1}\right|+\left|J L_{1}\right|_{G}\right)
$$

Proof. It is straightforward to compute

$$
\lambda-J_{\#} L_{\#}=\left(I-\left(J L_{1}+J_{1} L_{\#}\right)(\lambda-J L)^{-1}\right)(\lambda-J L) .
$$

According to assumption (A3), $J L_{1}(\lambda-J L)^{-1}$ is a closed operator with the domain $X$. The closed graph theorem implies that it is actually bounded with

$$
\begin{aligned}
\left|J L_{1}(\lambda-J L)^{-1}\right| & \leq\left|J L_{1}\right|_{G}\left(\left|(\lambda-J L)^{-1}\right|+\left|J L(\lambda-J L)^{-1}\right|\right) \\
& \leq\left|J L_{1}\right|_{G}\left(1+(1+|\lambda|)\left|(\lambda-J L)^{-1}\right|\right)
\end{aligned}
$$

where $J L=\lambda-(\lambda-J L)$ was used in the last step. The conclusion of the lemma follows from this along with the boundedness of $J_{1}, L$, and $L_{1}$.

### 9.1. Persistent exponential trichotomy and stability

In this section, our main task is to prove Theorem 2.4 as well as Proposition 2.9. With the help of Lemmas 9.2 and 6.2, we are able to prove most of Theorem 2.4 and Proposition 2.9 by standard arguments in the spectral theory. However, proving (2.26) requires more elaborated arguments as one of the perturbation term $J L_{1}$ is not necessarily a small bounded operator.
Proof of Theorem 2.4 except (2.26). Adopt the notation used in (2.24)

$$
\begin{equation*}
\epsilon \triangleq\left|J_{1}\right|+\left|L_{1}\right|+\left|J L_{1}\right|_{G} \tag{9.1}
\end{equation*}
$$

Let

$$
\sigma_{u, s}=\{\lambda \in \sigma(J L) \mid \pm \operatorname{Re} \lambda>0\} \subset \sigma(J L)
$$

and $\Omega_{u} \subset \mathbf{C}$ be open and bounded with smooth boundary $\Gamma_{u}=\partial \Omega_{u} \subset \mathbf{C} \backslash \sigma(J L)$ such that $\sigma(J L) \cap \Omega_{u}=\sigma_{u}$. According to Lemma 3.6, $\sigma_{s}$ is symmetric to $\sigma_{u}$ about $i \mathbf{R}$ and thus we let $\Omega_{s}$ be the domain symmetric to $\Omega_{u}$ and $\Gamma_{s}=\partial \Omega_{s}$. For small $\epsilon$, Lemma 9.2 allows us to define the following objects via standard contour integrals

$$
\begin{aligned}
\tilde{P}_{\#}^{u, s} & =\frac{1}{2 \pi i} \oint_{\Gamma_{u, s}}\left(z-J_{\#} L_{\#}\right)^{-1} d z, \quad E_{\#}^{u, s}=\tilde{P}_{\#}^{u, s} X \\
A_{\#}^{u, s} & =\left.\left(J_{\#} L_{\#}\right)\right|_{E_{\#}^{u, s}}=\frac{1}{2 \pi i} \oint_{\Gamma_{u, s}} z\left(z-J_{\#} L_{\#}\right)^{-1} d z
\end{aligned}
$$

Let

$$
\tilde{P}_{\#}^{c}=I-\tilde{P}_{\#}^{u}-\tilde{P}_{\#}^{s}, \quad E_{\#}^{c}=\tilde{P}_{\#}^{c} X, \quad A_{\#}^{c}=\left.\left(J_{\#} L_{\#}\right)\right|_{E_{\#}^{c}}
$$

and $\tilde{P}^{u, s}, E^{u, s, c}, A^{u, s, c}$ denote the corresponding unperturbed objects.
From the standard spectral theory, subspaces $E_{\#}^{u, s, c}$ are invariant under $J_{\#} L_{\#}$. Therefore, $A_{\#}^{u, s, c}$ are operators on $E_{\#}^{u, s, c}$ with $\sigma\left(A_{\#}^{u, s}\right) \subset \Omega_{u, s}$ and $\sigma\left(A_{\#}^{c}\right) \subset\left(\mathbf{C} \backslash\left(\Omega_{u} \cup\right.\right.$ $\left.\Omega_{s}\right)$ ). Lemma 9.2 implies

$$
\left|\tilde{P}_{\#}^{u, s}-\tilde{P}^{u, s}\right| \leq C \epsilon
$$

and thus $E_{\#}^{c}$ is $O(\epsilon)$ close to $E^{c}$, too. Along with the non-degeneracy of $\langle L \cdot, \cdot\rangle$ on $E^{u} \oplus E^{s}$ and $\left|L_{1}\right| \leq \epsilon$, above implies the non-degeneracy of $\left\langle L_{\#} \cdot, \cdot\right\rangle$ on $E_{\#}^{u} \oplus E_{\#}^{s}$. Therefore, we obtain from $X=E_{\#}^{u} \oplus E_{\#}^{s} \oplus E_{\#}^{c}$ and Lemma 6.2

$$
E_{\#}^{c}=\left\{u \in X \mid\left\langle L_{\#} u, v\right\rangle=0, \forall u \in E_{\#}^{u} \oplus E_{\#}^{s}\right\}
$$

As $O(\epsilon)$ perturbations, it is clear that subspaces $E_{\#}^{u, s, c}$ can be written as graphs of $O(\epsilon)$ bounded operators $S_{\#}^{u, s, c}$ in the coordinate frame $X=E^{u} \oplus E^{s} \oplus E^{c}$. Moreover, from the above integral forms, $A_{\#}^{u, s}$ are only $O(\epsilon)$ bounded perturbations to $J L$ on finite dimensional subspaces $E_{\#}^{u, s}$ which are $O(\epsilon)$ perturbations to $E^{u, s}$, and thus inequality (2.25) follows as well. Since the subspace $E_{\#}^{u} \oplus E_{\#}^{s}$, invariant under $J_{\#} L_{\#}$, is finite dimensional, the vanishness of $\left\langle L_{\#}, \cdot\right\rangle$ on $E_{\#}^{u, s}$ follows from Lemma 3.3. Through this point we complete the proof of parts (a) and (b), except (2.26), of Theorem 2.4.

Suppose, as in part (c) in Theorem 2.4, there exists $\delta>0$ such that $\langle L u, u\rangle \geq$ $\delta\|u\|^{2}$ for all $u \in E^{c}$. Since $L$ and $L_{1}$ are bounded and $E_{\#}^{c}$ is $O(\epsilon)$ perturbation of $E^{c}$, we have $\left\langle L_{\#} u, u\right\rangle \geq \frac{\delta}{2}\|u\|^{2}$ for all $u \in E_{\#}^{c}$. Therefore, the conservation of $\left\langle L_{\#} \cdot, \cdot\right\rangle$ by $e^{t J_{\#} L_{\#}}$ and the invariance of $E_{\#}^{c}$ under $e^{t J_{\#} L_{\#}}$ imply the boundedness of $\left.e^{t J_{\#} L_{\#}}\right|_{E_{\#}^{c}}$ uniformly in $t \in \mathbf{R}$ which proves part (b) of Theorem 2.4.

To complete the proof of Theorem 2.4, we shall prove the weak exponential growth estimate (2.26) in the perturbed center subspace $E_{\#}^{c}$, which involves much more than simple applications of the standard operator calculus and the conservation of energy. We first consider a special case where $J L$ has no hyperbolic directions.

Lemma 9.3. Assume $E^{c}=X$, then (2.26) holds for some $C, \epsilon_{0}>0$ depending on $J, L$.

Proof. From the construction of $E_{\#}^{u, s, c}$ and the additional assumption $E^{c}=$ $X$, it is clear $E_{\#}^{c}=X$ and (2.26) is reduced to

$$
\left|e^{t J_{\#} L_{\#}}\right| \leq C e^{C \epsilon|t|}, \quad \forall t \in \mathbf{R}, \quad \text { where } \epsilon \triangleq\left|J_{1}\right|+\left|L_{1}\right|+\left|J L_{1}\right|_{G}
$$

Since

$$
J_{\#} L_{\#}=J L_{\#}+J_{1}\left(L+L_{1}\right)
$$

and $J_{1}, L$, and $L_{1}$ are bounded with $\left|J_{1}\right| \leq \epsilon$, it suffices to prove

$$
\begin{equation*}
\left|e^{t J L_{\#}}\right| \leq C e^{C \epsilon|t|}, \quad \forall t \in \mathbf{R} \tag{9.2}
\end{equation*}
$$

Let $X=\oplus_{j=0}^{4} X_{j}$ be the decomposition, associated with projections $P_{j}$, given by Theorem 2.1 for $J$ and $L$, where $X_{0}=\operatorname{ker} L$ and $X_{5}=X_{6}=\{0\}$ due to the assumption $E^{c}=X$. Much as in (12.1), let

$$
L_{1, j k}=i_{j}^{*} L_{1} i_{k}: X_{k} \rightarrow X_{j}^{*}, \quad j, k=0, \ldots, 4
$$

which satisfy $L_{1, j k}=L_{1, k j}^{*}$ and

$$
L_{1}=\Sigma_{j, k=0}^{4} P_{j}^{*} L_{1, j k} P_{k}, \quad\left\langle L_{1, j k} u, v\right\rangle=\left\langle L_{1} u, v\right\rangle, \quad \forall u \in X_{k}, v \in X_{j}
$$

Let $J_{j k}=P_{j} J P_{k}^{*}$ be the blocks of $J$ associated to this decomposition, which have the forms given in Corollary 2.1 and satisfy

$$
J=\Sigma_{j, k=0}^{4} i_{X_{j}} J_{j k} i_{X_{k}}^{*}, \quad\left|J-i_{X_{3}} J_{33} i_{X_{3}}^{*}\right| \leq C
$$

We write

$$
\begin{aligned}
J L_{\#}= & J L+J L_{1} P_{3}+\Sigma_{k \in\{0,1,2,4\}} J L_{1} P_{k} \\
= & J L+\left(J-i_{X_{3}} J_{33} i_{X_{3}}^{*}\right) L_{1} P_{3} \\
& \quad+i_{X_{3}} J_{33} i_{X_{3}}^{*} \Sigma_{j, k=0}^{4} P_{j}^{*} L_{1, j k} P_{k} P_{3}+\Sigma_{k \in\{0,1,2,4\}} J L_{1} P_{k} \\
= & J L+i_{X_{3}} J_{33} L_{1,33} P_{3}+\left(J-i_{X_{3}} J_{33} i_{X_{3}}^{*}\right) L_{1} P_{3}+\Sigma_{k \in\{0,1,2,4\}} J L_{1} P_{k},
\end{aligned}
$$

where $P_{X_{j}} i_{X_{k}}=\delta_{j k} I_{X_{k}}$ is used. Since, for $k \neq 3, X_{k} \subset D(J L) \subset D\left(J L_{1}\right)$, we have that $J L_{1} P_{k}$ is a bounded operator with the norm bounded in terms of $\left|J L_{1}\right|_{G},\left|P_{k}\right|$, and $\left|J L P_{k}\right|$. Along with the boundedness of $J-i_{X_{3}} J_{33} i_{X_{3}}^{*}$, we obtain

$$
\begin{equation*}
\left|J L_{\#}-\left(J L+i_{X_{3}} J_{33} L_{1,33} P_{3}\right)\right|<C \epsilon \tag{9.3}
\end{equation*}
$$

From Theorem 2.1 we have

$$
J L+i_{X_{3}} J_{33} L_{1,33} P_{3} \longleftrightarrow\left(\begin{array}{ccccc}
0 & A_{01} & A_{02} & A_{03} & A_{04} \\
0 & A_{1} & A_{12} & A_{13} & A_{14} \\
0 & 0 & A_{2} & 0 & A_{24} \\
0 & 0 & 0 & A_{3}+J_{33} L_{1,33} & A_{34} \\
0 & 0 & 0 & 0 & A_{4}
\end{array}\right)
$$

Note all blocks of $J L+i_{X_{3}} J_{33} L_{1,33} P_{3}$ are identical to those of $J L$ except its (4,4)block

$$
A_{3}+J_{33} L_{1,33}=J_{33}\left(L_{X_{3}}+L_{1,33}\right)
$$

Since $\left\langle L_{X_{3}} \cdot, \cdot\right\rangle$ is uniformly positive on $X_{3}$, so is $\left\langle\left(L_{X_{3}}+L_{1,33}\right) \cdot, \cdot\right\rangle$. Therefore, the group $e^{t\left(A_{3}+J_{33} L_{1,33}\right)}$, conserving $\left\langle\left(L_{X_{3}}+L_{1,33}\right) \cdot, \cdot\right\rangle$, satisfies

$$
\left|e^{t\left(A_{3}+J_{33} L_{1,33}\right)}\right| \leq C, \quad \forall t \in \mathbf{R}
$$

By using the upper triangular form of $J L$, this inequality, assumption $E^{c}=X$ that $J L$ has no hyperbolic eigenvalues, and the finite dimensionality $\operatorname{dim} X_{1}=\operatorname{dim} X_{4}=$ $n^{-}(L)-\operatorname{dim} X_{2}$, it is easy to prove

$$
\left|e^{t\left(J L+i_{X_{3}} J_{33} L_{1,33} P_{3}\right)}\right| \leq C\left(1+|t|^{1+2 n^{-}(L)}\right), \quad \forall t \in \mathbf{R}
$$

Along with (9.3), above estimate implies (9.2) and we obtain (2.26) assuming $E^{c}=$ $X$.

By using the invariance of $E_{\#}^{c}$ under $J_{\#} L_{\#}$, in the following we convert $\left.e^{t J_{\#} L_{\#}}\right|_{E_{\#}^{c}}$ to a flow $e^{t \tilde{J}_{\#} \tilde{L}_{\#}}$ on $E^{c}$ via a similarity transformation and then apply Lemma 9.3 to obtain (2.26).

Proof of (2.26) in Theorem 2.4. In the general case, let $E_{\#}^{u, s, c}$ be the invariant unstable/stable/center subspaces and $P_{\#}^{u, s, c}$ be the projections associated to the decomposition $X=E_{\#}^{u} \oplus E_{\#}^{s} \oplus E_{\#}^{c}$. We also adopt the notations $E_{\#}^{u s}=E_{\#}^{u} \oplus E_{\#}^{s}$ and $P_{\#}^{u s}=P_{\#}^{u}+P_{\#}^{s}$. Correspondingly, let $E^{u, s, c, u s}$ and $P^{u, s, c, u s}$ denoted the unperturbed invariant subspaces and projections. Recall $E_{\#}^{c}$ can be written as the graph of a bounded operator $S_{\#}^{c}: E^{c} \rightarrow E^{u s}$ with $\left|S_{\#}^{c}\right|=O(\epsilon)$. Let $\tilde{S}_{\#}^{c}=i_{E^{c}}+S_{\#}^{c}$ : $E^{c} \rightarrow E_{\#}^{c} \subset X$ so that $E_{\#}^{c}=\tilde{S}_{\#}^{c}\left(X^{c}\right)$. Clearly, $P^{c} i_{E_{\#}^{c}}=\left(\tilde{S}_{\#}^{c}\right)^{-1}$.

Let

$$
J^{c}=P^{c} J\left(P^{c}\right)^{*}, \quad J_{\#}^{c}=P_{\#}^{c} J_{\#}\left(P_{\#}^{c}\right)^{*}, \quad L_{\#}^{c}=i_{E_{\#}^{c}}^{*} L_{\#} i_{E_{\#}^{c}}, \quad L^{c}=i_{E^{c}}^{*} L i_{E^{c}} .
$$

From the invariance of $E_{\#}^{c}$ under $J_{\#} L_{\#}$, the $L_{\#}$-orthogonality between $E_{\#}^{c}$ and $E_{\#}^{u s}$, and Lemma 12.3 applied to the decomposition $X=E_{\#}^{c} \oplus E_{\#}^{u s}$, we have

$$
J_{\#}^{c} L_{\#}^{c}=\left.J_{\#} L_{\#}\right|_{E_{\#}^{c}},\left.\quad e^{t\left(J_{\#} L_{\#}\right)}\right|_{E_{\#}^{c}}=e^{t J_{\#}^{c} L_{\#}^{c}}
$$

and the combination $\left(E_{\#}^{c}, J_{\#}^{c}, L_{\#}^{c}\right)$ satisfies (H1-3). Using the mapping $\tilde{S}_{\#}^{c}$, we may just consider its conjugate flow on $E^{c}$

$$
P^{c} i_{E_{\#}^{c}} e^{t J_{\#}^{c} L_{\#}^{c}} \tilde{S}_{\#}^{c}, \text { with the generator } P^{c} i_{E_{\#}^{c}} J_{\#}^{c} L_{\#}^{c} \tilde{S}_{\#}^{c} \text { on } E^{c}
$$

Let

$$
\tilde{L}_{\#}=\left(\tilde{S}_{\#}^{c}\right)^{*} L_{\#}^{c} \tilde{S}_{\#}^{c}=\left(\tilde{S}_{\#}^{c}\right)^{*} i_{E_{\#}^{c}}^{*} L_{\#} i_{E_{\#}^{c}} \tilde{S}_{\#}^{c}: E^{c} \rightarrow\left(E^{c}\right)^{*}
$$

and

$$
\tilde{J}_{\#}=P^{c} i_{E_{\#}^{c}} J_{\#}^{c}\left(P^{c} i_{E_{\#}^{c}}\right)^{*}=P^{c} P_{\#}^{c} J_{\#}\left(P^{c} P_{\#}^{c}\right)^{*}:\left(E^{c}\right)^{*} \supset D\left(\tilde{J}_{\#}\right) \rightarrow E^{c}
$$

Clearly,

$$
\begin{equation*}
\tilde{J}_{\#} \tilde{L}_{\#}=P^{c} i_{E_{\#}^{c}} J_{\#}^{c} L_{\#}^{c} \tilde{S}_{\#}^{c}=P^{c} i_{E_{\#}^{c}} J_{\#} L_{\#} \tilde{S}_{\#}^{c} \tag{9.4}
\end{equation*}
$$

Since $\left|P^{c} i_{E_{\#}^{c}}\right|,\left|\tilde{S}_{\#}^{c}\right| \leq 2$, in order to prove (2.26), it suffices to prove on $E^{c}$

$$
\begin{equation*}
\left|e^{\tilde{J}_{\#} \tilde{L}_{\#}}\right| \leq C e^{C \epsilon|t|}, \forall t \in \mathbf{R} \tag{9.5}
\end{equation*}
$$

for some $C$ depending only on $J$ and $L$. Our strategy is to verify that ( $E^{c}, \tilde{J}_{\#}, \tilde{L}_{\#}$ ) as a perturbation to $\left(E^{c}, J^{c}, L^{c}\right)$ satisfies (A1-3) and then apply Lemma 9.3.

When $\epsilon=0$, Lemma 12.3 ensures that the unperturbed

$$
\left(E^{c}, \tilde{J}_{\#}=J^{c}, \tilde{L}_{\#}=L^{c}=i_{E^{c}}^{*} L i_{E^{c}}\right)
$$

satisfies (H1-3). Moreover, since $\langle L \cdot, \cdot\rangle$ is non-degenerate on $E^{u s}$, we have
$\operatorname{dim} \operatorname{ker} L^{c}=\operatorname{dim} \operatorname{ker}\left(i_{E^{c}}^{*} L i_{E^{c}}\right)=\operatorname{dim} \operatorname{ker} L<\infty$
due to the $L$-orthogonality between $E^{c}$ and $E^{u s}$ and thus (A2) is satisfied by $\tilde{L}_{\#}$ for $\epsilon=0$. From the definitions, $\tilde{J}_{\#}-J^{c}$ is clearly anti-symmetric. We will show that it is also bounded. Using the fact $I-P_{\#}^{c}=P_{\#}^{u s}$, one may compute

$$
\begin{aligned}
\tilde{J}_{\#}-J^{c} & =-P^{c} P_{\#}^{u s} J_{\#}\left(P^{c} P_{\#}^{c}\right)^{*}-P^{c} J_{\#}\left(P^{c} P_{\#}^{u s}\right)^{*}+P^{c} J_{1}\left(P^{c}\right)^{*} \\
& =-P^{c} P_{\#}^{u s} P_{\#}^{u s} J_{\#}\left(P^{c} P_{\#}^{c}\right)^{*}-P^{c} J_{\#}\left(P_{\#}^{u s}\right)^{*}\left(P^{c} P_{\#}^{u s}\right)^{*}+P^{c} J_{1}\left(P^{c}\right)^{*}
\end{aligned}
$$

Due to the $L_{\# \text {-orthogonality between }} E_{\#}^{u s}$ and $E_{\#}^{c}$ and the non-degeneracy of $\left\langle L_{\#} \cdot, \cdot\right\rangle$ on $E_{\#}^{u s}$, it is straightforward to obtain that $L_{\#}$ is an isomorphism from $E_{\#}^{u s}$ to $R\left(\left(P_{\#}^{u s}\right)^{*}\right)=\left(P_{\#}^{u s}\right)^{*}\left(E_{\#}^{u s}\right)^{*}$. Since $E_{\#}^{u s} \subset D\left(J_{\#} L_{\#}\right)$, we have $R\left(\left(P_{\#}^{u s}\right)^{*}\right) \subset$ $D\left(J_{\#}\right)$ and thus $\left.J_{\#}\right|_{R\left(\left(P_{\#}^{u s}\right)^{*}\right)}$ is a bounded operator. To estimate its norm, we use the relationship

$$
\left.J_{\#}\right|_{R\left(\left(P_{\#}^{u s}\right)^{*}\right)}=\left(\left.J_{\#} L_{\#}\right|_{E_{\#}^{u s}}\right)\left(\left.L_{\#}\right|_{E_{\#}^{u s}}\right)^{-1}=\left(A_{\#}^{u} \oplus A_{\#}^{s}\right)\left(\left.L_{\#}\right|_{E_{\#}^{u s}}\right)^{-1} .
$$

Recall $E_{\#}^{u s}$ is $O(\epsilon)$ perturbation to $E^{u s}$ and $L_{\#}$ is $O(\epsilon)$ to $L$. Moreover, the spectral integral representations of $A_{\#}^{u, s}$ yield that they are $O(\epsilon)$ perturbation to $\left.J L\right|_{E^{u s}}$. Therefore, we obtain that

$$
\left|J_{\#}\right|_{R\left(\left(P_{\#}^{u s}\right)^{*}\right.}|\leq C \Longrightarrow| P_{\#}^{u s} J_{\#}\left|=\left|J_{\#}\left(P_{\#}^{u s}\right)^{*}\right| \leq C\right.
$$

for some $C>0$ depending on $J$ and $L$. Since $\left|P^{c} P_{\#}^{u s}\right| \leq C \epsilon$, we have

$$
\left|\tilde{J}_{1, \#}\right| \leq C \epsilon, \quad \text { where } \tilde{J}_{1, \#} \triangleq \tilde{J}_{\#}-J^{c}
$$

From the definition of $\tilde{L}_{\#}$, it is easy to obtain

$$
\tilde{L}_{1, \#}=\tilde{L}_{1, \#}^{*}, \quad\left|\tilde{L}_{1, \#}\right| \leq C \epsilon, \quad \text { where } \tilde{L}_{1, \#} \triangleq \tilde{L}_{\#}-L^{c}
$$

Therefore, we finish verifying (A1) for $\left(E^{c}, \tilde{J}_{\#}, \tilde{L}_{\#}\right)$.
We proceed to verify (A3). From Lemma 9.2, we have $D\left(J_{\#} L_{\#}\right)=D(J L)$. Since $E_{\#}^{c}=\tilde{S}_{\#}^{c}\left(E^{c}\right)$ is the graph of $S_{\#}^{c}: E^{c} \rightarrow E^{u s}$ and $E^{u s} \subset D(J L)=D\left(J_{\#} L_{\#}\right)$, we obtain

$$
D\left(J_{\#} L_{\#}\right) \cap E_{\#}^{c}=\tilde{S}_{\#}^{c}\left(E^{c} \cap D(J L)\right)
$$

From the boundedness of $\tilde{J}_{1, \#}$ and (9.4), we further obtain

$$
D\left(J^{c} \tilde{L}_{\#}\right)=D\left(\tilde{J}_{\#} \tilde{L}_{\#}\right)=E^{c} \cap D(J L)=D\left(J^{c} L^{c}\right)
$$

which along with $\tilde{L}_{\#}=L^{c}+\tilde{L}_{1, \#}$ obviously implies $D\left(J^{c} L^{c}\right) \subset D\left(J^{c} \tilde{L}_{1, \#}\right)$.
In the next we estimate the graph norm of $J^{c} \tilde{L}_{1, \#}$, like the one defined in (2.23), on the domain $D\left(J^{c} \tilde{L}_{\#}\right)=E^{c} \cap D(J L)$. From (9.4) one may compute that, when restricted on $E^{c} \cap D(J L)$,

$$
\begin{equation*}
\tilde{J}_{\#} \tilde{L}_{\#}-J^{c} L^{c}=-P^{u s} i_{E_{\#}^{c}} J_{\#} L_{\#} \tilde{S}_{\#}^{c}+J_{\#} L_{\#} S_{\#}^{c}+J_{\#} L_{\#}-J L \tag{9.6}
\end{equation*}
$$

We shall use

$$
\begin{equation*}
J_{\#} L_{\#}=J L+J L_{1}+J_{1} L_{\#} \tag{9.7}
\end{equation*}
$$

to estimate the three terms in (9.6). In fact, for any $v \in D(J L)$,

$$
\left\|J_{\#} L_{\#} v-J L v\right\| \leq\left|J L_{1}\right|_{G}\|v\|_{G}+\mid J_{1} L_{\#}\| \| v\|\leq C \epsilon\| v \|_{G}
$$

for some $C>0$ depending on $J$ and $L$. A combination of this inequality with (9.6) and the fact $\left|P^{u s} i_{E_{\#}^{c}}\right| \leq C \epsilon$ implies, for any $u \in E^{c} \cap D(J L)$,

$$
\left\|\left(\tilde{J}_{\#} \tilde{L}_{\#}-J^{c} L^{c}\right) u\right\| \leq C\left(\epsilon\left\|\tilde{S}_{\#}^{c} u\right\|_{G}+\left\|S_{\#}^{c} u\right\|_{G}+\epsilon\|u\|_{G}\right) \leq C \epsilon\|u\|_{G}
$$

where we used the fact that $J L$ is bounded on $E^{u s} \subset D(J L)$. Since

$$
J^{c} \tilde{L}_{1, \#}=\tilde{J}_{\#} \tilde{L}_{\#}-J^{c} L^{c}-\tilde{J}_{1, \#} \tilde{L}_{\#}
$$

with $\left|\tilde{J}_{1, \#}\right| \leq C \epsilon$, the above inequality implies

$$
\left\|J^{c} \tilde{L}_{1, \#} u\right\| \leq C \epsilon\|u\|_{G}, \quad \forall u \in E^{c} \cap D(J L)
$$

for some $C>0$ depending on $J$ and $L$. The above estimates allow us to apply Lemma 9.3 to obtain (9.5) for $e^{t \tilde{J}_{\#} \tilde{L}_{\#}}$ on $E^{c}$, which in turn implies (2.26) for $e^{t J_{\#} L_{\#}}$ on $E_{\#}^{c}$.

To complete this section we present
Proof of Proposition 2.9. Adopt the notations used in (2.24)-(9.1), and let

$$
\epsilon \triangleq\left|J_{1}\right|+\left|L_{1}\right|+\left|J L_{1}\right|_{G}
$$

Let $\Omega \subset \mathbf{C}$ be an open domain with the compact closure and smooth boundary $\Gamma \subset \mathbf{C} \backslash i \mathbf{R}$ such that $\Omega \cap \sigma(J L)=\sigma_{2}$. For small $\epsilon$, Lemma 9.2 allows us to define the following objects via standard contour integrals

$$
\begin{aligned}
& P_{\#}=\frac{1}{2 \pi i} \oint_{\Gamma}\left(z-J_{\#} L_{\#}\right)^{-1} d z, \quad X_{2}^{\#}=P_{\#} X \subset D\left(J_{\#} L_{\#}\right), \quad X_{1}^{\#}=\left(I-P_{\#}\right) X \\
& A_{1,2}^{\#}=\left.\left(J_{\#} L_{\#}\right)\right|_{X_{1,2}}=\frac{1}{2 \pi i} \oint_{\Gamma} z\left(z-J_{\#} L_{\#}\right)^{-1} d z
\end{aligned}
$$

Let $P, X_{1,2}$, and $A_{1,2}$ denote the corresponding unperturbed objects.
From the standard spectral theory, the decomposition $X=X_{1}^{\#} \oplus X_{2}^{\#}$ is invariant under $J_{\#} L_{\#}$ and thus $A_{1,2}^{\#}$ are operators on $X_{1,2}^{\#}$ with $\sigma\left(A_{2}^{\#}\right) \subset \Omega$. (In fact $\left.\sigma\left(A_{1,2}\right)=\sigma_{1,2}.\right)$ Since $\sigma\left(J_{\#} L_{\#}\right)=\sigma\left(A_{1}^{\#}\right) \cup \sigma\left(A_{2}^{\#}\right)$, we only need to prove that $\sigma\left(A_{1}^{\#}\right) \subset i \mathbf{R}$.

Since the decomposition $X=X_{1} \oplus X_{2}$ is $L$-orthogonal and $X_{2} \subset D(J L)$, Lemma 12.3 implies that $\left(X_{1}, J_{X_{1}}=(I-P) J(I-P)^{*}, L_{X_{1}}\right)$ satisfies (H1-3). Therefore, the index theorem Theorem 2.3 applies to $L_{X_{1}}$ and $\left.J L\right|_{X_{1}}$ which along with the second assumption of Proposition 2.9 implies that $n^{-}\left(L_{X_{1}}\right)=0$. As $L_{X_{1}}$ satisfies (H2), we obtain that $L_{X_{1}}$ is positive definite. Lemma 9.2 implies

$$
\left|P_{\#}-P\right| \leq C \epsilon
$$

and thus $X_{1}^{\#}$ is $O(\epsilon)$ close to $X_{1}$. Namely $X_{1}^{\#}$ can be written as the graph of an $O(\epsilon)$ order bounded operator $S_{\#}: X_{1} \rightarrow X_{2}$. It immediately implies that $L_{X_{1}^{\#}}$ is uniformly positive on $X_{1}^{\#}$ and the proposition follows from the invariance of $X_{1}^{\#}$ under $J_{\#} L_{\#}$.

### 9.2. Perturbations of purely imaginary spectra

In this section, we consider $\sigma\left(J_{\#} L_{\#}\right)$ near some $i \mu \in \sigma(J L) \cap i \mathbf{R}$ and prove Theorems 2.5 and 2.6.
'Structurally stable' cases. We still adopt the notation used in (2.24) and let

$$
\begin{equation*}
\epsilon \triangleq\left|J_{1}\right|+\left|L_{1}\right|+\left|J L_{1}\right|_{G} \tag{9.8}
\end{equation*}
$$

Case 1: $i \mu \in \sigma(J L) \cap i \mathbf{R}$ is isolated with $\langle L \cdot, \cdot\rangle$ sign definite on $E_{i \mu}$.
Suppose $\delta>0$ and $\langle L u, u\rangle \geq \delta\|u\|^{2}$, for all $u \in E_{i \mu}$ (the opposite case $\langle L \cdot, \cdot\rangle \leq$ $-\delta<0$ on $E_{i \mu}$ is similar). Since $i \mu$ is assumed to be isolated in $\sigma(J L)$, there exists $\alpha>0$ such that the closed disk $\overline{B(i \mu, \alpha)} \cap \sigma(J L)=\{i \mu\}$. Let $\Gamma=\partial B(i \mu, \alpha)$ and $\Gamma \cap \sigma\left(J_{\#} L_{\#}\right)=\emptyset$ for small $\epsilon$ due to Lemma 9.2. Define

$$
\tilde{P}_{\#}=\frac{1}{2 \pi i} \oint_{\Gamma}\left(z-J_{\#} L_{\#}\right)^{-1} d z, \quad \tilde{E}_{\#}=\tilde{P}_{\#} X
$$

From the standard spectral theory, $\tilde{E}_{\#}$ is invariant under $J_{\#} L_{\#}$ and

$$
\sigma\left(J_{\#} L_{\#}\right) \cap \overline{B(i \mu, \alpha)}=\sigma\left(\left.J_{\#} L_{\#}\right|_{\tilde{E}_{\#}}\right)
$$

Lemma 9.2, the isolation of $i \mu$, and Proposition 2.3 imply that $\tilde{E}_{\#}$ is $O(\epsilon)$ close to $E_{i \mu}$. The positive definiteness assumption of $\langle L \cdot, \cdot\rangle$ on $E_{i \mu}$ and the boundedness of $L$ and $L_{1}$ imply that $\left\langle L_{\#} \cdot, \cdot\right\rangle$ is also positive definite on $\tilde{E}_{\#}$. The stability both forward and backward in time - of $e^{t J_{\#} L_{\#}}$ on $\tilde{E}_{\#}$, due to the conservation of energy, implies

$$
\sigma\left(J_{\#} L_{\#}\right) \cap \overline{B(i \mu, \alpha)}=\sigma\left(\left.J_{\#} L_{\#}\right|_{\tilde{E}_{\#}}\right) \subset i \mathbf{R}
$$

Case 2: $i \mu \in \sigma(J L) \cap i \mathbf{R}$ and $\langle L \cdot, \cdot\rangle$ is positive definite on $E_{i \mu}$.
Then

$$
\begin{equation*}
E_{i \mu}=\{0\} \text { or }\langle L u, u\rangle \geq \delta\|u\|^{2}, \forall u \in E_{i \mu} \tag{9.9}
\end{equation*}
$$

for some $\delta>0$.
Remark 9.1. In this case, besides the possibility of an isolated eigenvalue $i \mu \in$ $\sigma(J L)$ with $L$ positive definite on $E_{i \mu}$, we are mainly concerned with the scenario that i $\mu$ is embedded in the continuous spectrum, whether an eigenvalue or not, but without any eigenvector in a non-positive direction of L. Our conclusion is that, under small perturbations, no hyperbolic eigenvalues (i.e. away from imaginary axis) may bifurcate from $i \mu$.

We argue by contradiction for Case 2. Suppose Theorem 2.5 does not hold in this case, then there exist a sequence

$$
J_{\# n}=J+\tilde{J}_{n}, L_{\# n}=L+\tilde{L}_{n}, \quad n=1,2, \ldots
$$

satisfying (A1-3) for each $n$ such that

$$
\exists \lambda_{n} \in \sigma\left(J_{\# n} L_{\# n}\right) \backslash i \mathbf{R} ; \quad \epsilon_{n} \triangleq\left|\tilde{J}_{n}\right|+\left|\tilde{L}_{n}\right|+\left|J \tilde{L}_{n}\right|_{G} \rightarrow 0 ; \quad \delta_{n} \triangleq\left|\lambda_{n}-i \mu\right| \rightarrow 0
$$

Since not in $i \mathbf{R}, \lambda_{n}$ must be eigenvalues. Let

$$
u_{n} \in X, \quad J_{\# n} L_{\# n} u_{n}=\lambda_{n} u_{n}, \quad\left\|u_{n}\right\|=1
$$

Using the graph norm of $J L_{1}$, one may estimate

$$
\left\|J \tilde{L}_{n} u_{n}\right\| \leq\left|J \tilde{L}_{n}\right|_{G}\left(1+\left\|J L u_{n}\right\|\right) \leq\left|J \tilde{L}_{n}\right|_{G}\left(1+\left|\lambda_{n}\right|+\left\|J L u_{n}-\lambda_{n} u_{n}\right\|\right)
$$

and

$$
\left\|J L u_{n}-\lambda_{n} u_{n}\right\|=\left\|J L u_{n}-J_{\# n} L_{\# n} u_{n}\right\| \leq\left\|\tilde{J}_{n} L_{\# n} u_{n}\right\|+\left\|J \tilde{L}_{n} u_{n}\right\|
$$

Therefore, we obtain

$$
\begin{equation*}
\left\|J L u_{n}-\lambda_{n} u_{n}\right\| \leq C \epsilon_{n}, \quad\left\|J L u_{n}-i \mu u_{n}\right\| \leq\left(C \epsilon_{n}+\delta_{n}\right) \tag{9.10}
\end{equation*}
$$

for some $C>0$ depending on $|L|$ and $\mu$.
Let $X=\oplus_{j=0}^{6}$ be the decomposition given by Theorem 2.1 for $(L, J)$, with $X_{0}=\operatorname{ker} L, P_{j}$ be the associated projections, and $u_{n, j}=P_{j} u_{n}$. Let $A_{j}$ and $A_{j k}$ denote the blocks of $J L$ in this decomposition as given in Theorem 2.1. From the commutativity between $J L$ and $P_{5,6}$, we obtain from (9.10)

$$
\left\|A_{5} u_{n, 5}-i \mu u_{n, 5}\right\|+\left\|A_{6} u_{n, 6}-i \mu u_{n, 6}\right\| \leq C\left(\epsilon_{n}+\delta_{n}\right)
$$

Since $\sigma\left(A_{5,6}\right) \cap i \mathbf{R}=\emptyset$, we have

$$
\begin{equation*}
\left\|u_{n, 5}\right\|+\left\|u_{n, 6}\right\| \leq C\left(\epsilon_{n}+\delta_{n}\right) \tag{9.11}
\end{equation*}
$$

From Lemma 3.3, we have $\left\langle L_{\# n} u_{n}, u_{n}\right\rangle=0$. Along with (9.11) this implies that

$$
\begin{equation*}
\left|2\left\langle L u_{n, 1}, u_{n, 4}\right\rangle+\left\langle L_{2} u_{n, 2}, u_{n, 2}\right\rangle+\left\langle L_{3} u_{n, 3}, u_{n, 3}\right\rangle\right| \leq C\left(\epsilon_{n}+\delta_{n}\right) \tag{9.12}
\end{equation*}
$$

Applying $P_{j}, j=0, \ldots, 4$ to (9.10) and using Theorem 2.1, we have

$$
\begin{align*}
& \left\|A_{4} u_{n, 4}-i \mu u_{n, 4}\right\| \leq C\left(\epsilon_{n}+\delta_{n}\right)  \tag{9.13}\\
& \left\|A_{3} u_{n, 3}+A_{34} u_{n, 4}-i \mu u_{n, 3}\right\| \leq C\left(\epsilon_{n}+\delta_{n}\right) \\
& \left\|A_{2} u_{n, 2}+A_{24} u_{n, 4}-i \mu u_{n, 2}\right\| \leq C\left(\epsilon_{n}+\delta_{n}\right) \\
& \left\|A_{1} u_{n, 1}+A_{12} u_{n, 2}+A_{13} u_{n, 3}+A_{14} u_{n, 4}-i \mu u_{n, 1}\right\| \leq C\left(\epsilon_{n}+\delta_{n}\right) \\
& \left\|A_{01} u_{n, 1}+A_{02} u_{n, 2}+A_{03} u_{n, 3}+A_{04} u_{n, 4}-i \mu u_{n, 0}\right\| \leq C\left(\epsilon_{n}+\delta_{n}\right) \tag{9.14}
\end{align*}
$$

Since $\operatorname{dim} X_{j}<\infty$ when $j \neq 0,3$, subject to a subsequence, we may assume that as $n \rightarrow \infty$,

$$
u_{n, j} \rightarrow u_{j}, j=1,2,4 ; \quad u_{n, 5}, u_{n, 6} \rightarrow 0 ; \quad u_{n, j} \rightharpoonup u_{j}, j=0,3
$$

Passing to the limits in the above inequalities and using the boundedness of $A_{j}$ and $A_{j k}$ except $A_{3}$, we obtain

$$
\begin{aligned}
& A_{4} u_{4}-i \mu u_{4}=0 ; \quad A_{2} u_{2}+A_{24} u_{4}-i \mu u_{2}=0 \\
& A_{1} u_{1}+A_{12} u_{2}+A_{13} u_{3}+A_{14} u_{4}-i \mu u_{1}=0 \\
& A_{01} u_{1}+A_{02} u_{2}+A_{03} u_{3}+A_{04} u_{4}-i \mu u_{0}=0
\end{aligned}
$$

Moreover, the above inequality involving $A_{3} u_{n, 3}$ also implies that

$$
u_{n, 3} \rightharpoonup u_{3}, \quad A_{3} u_{n, 3} \rightharpoonup-A_{34} u_{4}+i \mu u_{3}
$$

Since the graph of the closed operator of $A_{3}$ as a closed subspace in $X_{3} \times X_{3}$ is also closed under the weak topology, we obtain

$$
u_{3} \in D\left(A_{3}\right) \text { and } A_{3} u_{3}+A_{34} u_{4}-i \mu u_{3}=0
$$

These equalities imply that

$$
J L u=i \mu u, \quad \text { where } u=u_{0}+u_{1}+u_{2}+u_{3}+u_{4}
$$

In addition, (9.12) implies

$$
\langle L u, u\rangle=2\left\langle L u_{1}, u_{4}\right\rangle+\left\langle L_{2} u_{2}, u_{2}\right\rangle+\left\langle L_{3} u_{3}, u_{3}\right\rangle \leq 0 .
$$

Due to property (9.9) of $i \mu$, we must have $u=0$, which immediately yields $u_{n, j} \rightarrow 0, j=1,2,4,5,6$ and thus (9.12) implies $u_{n, 3} \rightarrow 0$ as well. Then the normalization $\left\|u_{n}\right\|=1$ implies that we must have $\operatorname{dim} \operatorname{ker} L \geq 1,\left\|u_{n, 0}\right\| \rightarrow 1$ and $u_{n, 0} \rightharpoonup 0$. From (9.14), we obtain $\mu=0$. As ker $L$ is nontrivial, this again contradicts to (9.9). Therefore, Theorem 2.5 holds in this case.

Summarizing the above two cases, Theorem 2.5 is proved.
'Structurally unstable' cases. In the following, we will consider cases in Theorem 2.6 for the structural instability. In many applications the symplectic structure $J$ usually does not vary, therefore we will fix $J$ and focus on constructing perturbations to the energy operator $L$ to induce instabilities arising from a purely imaginary eigenvalue $i \mu$ of $J L$. Recall we have to complexify $X, J$, and $L$ accordingly. However, keep in mind that we would like to construct real perturbations to create unstable eigenvalues near $i \mu$. This would require the perturbations to also satisfy (12.12), see Remark 12.5. Recall that while $J L$ is a linear operator, $L$ and $J$ are
complexified as Hermitian forms or anti-linear mappings, see (12.8) and (12.9).
Case 3: $i \mu \in \sigma(J L)$ and $\exists$ a closed subspace $\{0\} \neq Y \subset E_{i \mu}$ such that

$$
J L(Y) \subset Y, \text { and }\langle L \cdot, \cdot\rangle \text { is non-degenerate and sign indefinite on } Y .
$$

REmARK 9.2. Clearly, this includes, but not limited to, the situation where $\langle L \cdot, \cdot\rangle$ is non-degenerate and indefinite on $E_{i \mu}$, a special case of which is when $i \mu$ is isolated in $\sigma(J L)$. It is analyzed in several subcases below.

We will construct a perturbation $L_{\#}$ such that $\sigma\left(J L_{\#}\right)$ contains a hyperbolic eigenvalue near $i \mu$. The proof will basically be carried out in some finite dimensional subspaces. Such finite dimensional problems had been well studied in the literature, mostly for the Case 3b below when there are two eigenvectors of opposite signs of $\langle L \cdot, \cdot\rangle$ (see e.g. [59, 24]). We could not find a reference for the proof of structural instability when the indefiniteness of $\left.L\right|_{E_{i \mu}}$ is caused by a Jordan chain of $J L$ (Case 3c below). So we give a detailed proof for the general case, which will also be used in later cases of embedded eigenvalues. Our proof for the Case 3c uses the special basis constructed in Proposition 2.2 for the Jordan blocks of $J L$ on $E_{i \mu}$.

Recall $\bar{u} \in E_{-i \mu}$ for any $u \in E_{i \mu}$. Let

$$
Y_{\mu}=\{u+\bar{v} \mid u, v \in Y\} \subset E_{i \mu}+E_{-i \mu}
$$

From Lemma 3.3 and the assumption on $Y,\langle L \cdot, \cdot\rangle$ is still non-degenerate on $Y_{\mu}$ which is also clearly invariant under $J L$. Recall that

$$
Y_{\mu}^{\perp_{L}}=\left\{u \in X \mid\langle L u, v\rangle=0, \forall v \in Y_{\mu}\right\} .
$$

From Lemmas 12.2 and 3.2, $Y_{\mu}^{\perp_{L}}$ is also invariant under $e^{t J L}$ and $X=Y_{\mu} \oplus Y_{\mu}^{\perp_{L}}$. The definition of $Y_{\mu}$ implies that $Y_{\mu}$ is real, in the sense $\bar{u} \in Y_{\mu}$ for any $u \in Y_{\mu}$, and thus the complexification of a real subspace of $X$. According to Lemma 12.3, $\left.J L\right|_{Y_{\mu}}$ is a also a Hamiltonian operator satisfying hypotheses (H1-3) with the nondegenerate energy $L_{Y_{\mu}}$, defined in (12.1). Therefore, we may apply Proposition 2.2 to $Y_{\mu}$ and $\left.J L\right|_{Y_{\mu}}$, where $Y \subset Y_{\mu}$ is the subspace of all generalized eigenvectors of $i \mu$ of $\left.J L\right|_{Y_{\mu}}$. Since $L$ is non-degenerate on $Y$, it is clear that Case $\mathbf{3}$ contains the following three subcases only.

Case 3a: $\mu \neq 0$ and $\langle L \cdot, \cdot\rangle$ changes sign on $\operatorname{ker}(J L-i \mu) \cap Y$.
In this subcase, let $u_{ \pm} \in \operatorname{ker}(J L-i \mu) \cap Y$ be such that $\pm\left\langle L u_{ \pm}, u_{ \pm}\right\rangle>0$. By a Gram-Schmidt process, without loss of generality, we may assume

$$
\left\langle L u_{ \pm}, u_{ \pm}\right\rangle= \pm 1, \quad\left\langle L u_{+}, u_{-}\right\rangle=0
$$

Note that $u_{ \pm}$can not be real for $\mu \neq 0$. As we will construct real perturbations to create instability, we have to consider the complex conjugate of $u_{ \pm}$as well. Let

$$
X_{1}=\operatorname{span}\left\{u_{+}, u_{-}, \overline{u_{+}}, \overline{u_{-}}\right\}, \quad X_{2}=X_{1}^{\perp_{L}}=\left\{v \in X \mid\left\langle L u_{ \pm}, v\right\rangle=\left\langle L \overline{u_{ \pm}}, v\right\rangle=0\right\}
$$

It is clear that subspaces $X_{1,2}$ are comlexifications of real subspaces in the sense

$$
\begin{equation*}
\bar{u} \in X_{1,2} \text { if } u \in X_{1,2} \tag{9.15}
\end{equation*}
$$

Note that $\overline{u_{ \pm}}$are eigenvectors of $-i \mu(\neq i \mu)$ and that $E_{i \mu}$ and $E_{-i \mu}$ are $L$ orthogonal. Therefore, from the complexification process, it is easy to verify that,
with respect to this basis of the invariant subspace $X_{1}$ of $J L$, operators $L_{X_{1}}$ and $\left.J L\right|_{X_{1}}$ take the forms

$$
L_{X_{1}}=\left(\begin{array}{cc}
\Lambda & 0 \\
0 & \Lambda
\end{array}\right),\left.\quad J L\right|_{X_{1}} \triangleq A_{1}=i \mu\left(\begin{array}{cc}
I_{2 \times 2} & 0 \\
0 & -I_{2 \times 2}
\end{array}\right)
$$

where

$$
\Lambda=\left(\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right), \quad I_{2 \times 2}=\left(\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right)
$$

From the invariance of $X_{1}$ and Lemmas 12.2 and 3.2, $X_{2}=X_{1}^{\perp_{L}}$ is also invariant under $J L$ and $X=X_{1} \oplus X_{2}$. In this decomposition, $L, J L$, and $J$ take the forms

$$
L=\left(\begin{array}{cc}
L_{X_{1}} & 0  \tag{9.16}\\
0 & L_{X_{2}}
\end{array}\right), \quad J L=\left(\begin{array}{cc}
A_{X_{1}} & 0 \\
0 & A_{X_{2}}
\end{array}\right), \quad J=\left(\begin{array}{cc}
J_{X_{1}} & 0 \\
0 & J_{X_{2}}
\end{array}\right)
$$

where, with respect to the basis of $X_{1}^{*}$ dual to $\left\{u_{ \pm}, \overline{u_{ \pm}}\right\}$,

$$
J_{X_{1}}=i \mu\left(\begin{array}{cc}
\Lambda & 0 \\
0 & -\Lambda
\end{array}\right), \quad J_{X_{2}}: X_{2}^{*} \supset D\left(J_{X_{2}}\right) \rightarrow X_{2}, \quad J_{X_{2}}^{*}=-J_{X_{2}}
$$

Here, $J^{*}=-J$ is used.
Consider a perturbation $L_{1}$ in the form of

$$
L_{1}=\left(\begin{array}{cc}
L_{1, X_{1}} & 0 \\
0 & 0
\end{array}\right), \quad \text { where } L_{1, X_{1}}=\left(\begin{array}{cc}
\epsilon R & 0 \\
0 & \epsilon R
\end{array}\right), \quad R=\left(\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right)
$$

It is straightforward to verify that $L_{1}$ is real, namely, $\left\langle L_{1} \bar{u}, \bar{v}\right\rangle=\overline{\left\langle L_{1} u, v\right\rangle}$. Let $L_{\#}=L+L_{1}$. Clearly, the decomposition $X=X_{1} \oplus X_{2}$ is still invariant under $J L_{\#}$ and orthogonal with respect to $L_{\#}$. Therefore, by a direct computation on the $4 \times 4$ matrix $\left.J L_{\#}\right|_{X_{1}}$ which can be further reduced to the $2 \times 2$ matrix $i \mu \Lambda(\Lambda+\epsilon R)$, we obtain

$$
i \mu \pm \epsilon \mu \in \sigma\left(J L_{\#}\right)
$$

Therefore, $\sigma\left(J L_{\#}\right)$ contains hyperbolic eigenvalues near $i \mu$ for any $\epsilon \neq 0$.
Case 3b: $\mu=0$ and $\langle L \cdot, \cdot\rangle$ changes sign on $\operatorname{ker}(J L-i \mu) \cap Y$.
In this case one may proceed as in the above through (9.16), however, with $J_{X_{1}}=0$. Therefore, no hyperbolic eigenvalue can bifurcate through such type of perturbations of $L$.

Cases 3c: $\mu \neq 0$ and $Y$ contains a non-trivial Jordan chain $u_{j}=(J L-i \mu)^{j-1} u_{1}$, $j=1, \ldots, k>1$, of $J L$ such that $u_{k} \in \operatorname{ker}(J L-i \mu) \backslash\{0\}$ and $\langle L \cdot, \cdot\rangle$ is nondegenerate on $\operatorname{span}\left\{u_{1}, \ldots, u_{k}\right\}$.

Again in this case let

$$
X_{1}=\operatorname{span}\left\{u_{j}, \overline{u_{j}} \mid j=1, \ldots, k\right\}, \quad X_{2}=X_{1}^{\perp_{L}}
$$

which is a $L$-orthogonal invariant decomposition under $J L$ satisfying (9.15) and thus the forms (9.16) hold. From Proposition 2.2, without loss of generality, we may assume that, with respect to the basis $\left\{u_{1}, \ldots, u_{k}, \overline{u_{1}}, \ldots, \overline{u_{k}}\right\}$ (as well as its dual basis in $\left.X_{1}^{*}\right), L_{X_{1}},\left.J L\right|_{X_{1}} \triangleq A_{X_{1}}$, and $J_{X_{1}}$ take the forms

$$
L_{X_{1}}=\left(\begin{array}{cc}
B_{+} & 0 \\
0 & \overline{B_{+}}
\end{array}\right), \quad A_{X_{1}}=\left(\begin{array}{cc}
A_{+} & 0 \\
0 & \frac{A_{+}}{+}
\end{array}\right), \quad J_{X_{1}}=\left(\begin{array}{cc}
J_{+} & 0 \\
0 & \frac{J_{+}}{}
\end{array}\right)
$$

where

$$
B_{+}=\left(\begin{array}{cccc}
0 & \ldots & 0 & b_{k} \\
0 & \ldots & b_{k-1} & 0 \\
\ldots & & & \\
b_{1} & \ldots & 0 & 0
\end{array}\right), \quad A_{+}=\left(\begin{array}{ccccc}
i \mu & 0 & \ldots & 0 & 0 \\
1 & i \mu & \ldots & 0 & 0 \\
\ldots & & & & \\
0 & 0 & \ldots & 1 & i \mu
\end{array}\right)
$$

and $b_{j+1}=-b_{j}, b_{k+1-j}=\overline{b_{j}}$. Therefore, $b_{j} \in\{ \pm i\}$ if $2 \mid k$ or $b_{j} \in\{ \pm 1\}$ otherwise. Then one may compute

$$
J_{+}=\left(\begin{array}{ccccc}
0 & 0 & \ldots & 0 & i \mu{\overline{b_{1}}}^{-1} \\
0 & 0 & \ldots & i \mu{\overline{b_{2}}}^{-1} & {\overline{b_{1}}}^{-1} \\
\cdots & & & & \\
i \mu{\overline{b_{k}}}^{-1} & {\overline{b_{k-1}}}^{-1} & \ldots & 0 & 0
\end{array}\right)
$$

Here, note that ${\overline{b_{j}}}^{-1}$ instead of $b_{j}^{-1}$ appears in above $J_{+}$, namely

$$
J_{+}\left(u_{j}^{*}\right)=i \mu{\overline{b_{k+1-j}}}^{-1} u_{k+1-j}+{\overline{b_{k+1-j}}}^{-1} u_{k+2-j}
$$

where $u_{k+1}=0$ is understood. This is due to the anti-linear complexification of $L$ and $J$, see (12.9). In fact, let $\left\{u_{j}^{*},{\overline{u_{j}}}^{*} \mid j=1, \ldots, k\right\}$ be the dual basis in $X_{1}^{*}$ which are complex linear functionals. We have

$$
\begin{aligned}
\left\langle u_{l}^{*}, J u_{j}^{*}\right\rangle & \left.=\left\langle u_{l}^{*}, J\left(b_{k+1-j}^{-1} L u_{k+1-j}\right)\right\rangle=\left\langle u_{l}^{*},{\overline{b_{k+1-j}}}^{-1} J L u_{k+1-j}\right)\right\rangle \\
& ={\overline{b_{k+1-j}}}^{-1}\left\langle u_{l}^{*}, J L u_{k+1-j}\right\rangle={\overline{b_{k+1-j}}}^{-1}\left\langle u_{l}^{*}, i \mu u_{k+1-j}+u_{k+2-j}\right\rangle \\
& ={\overline{b_{k+1-j}}}^{-1}\left(i \mu \delta_{l, k+1-j}+\delta_{l, k+2-j}\right) .
\end{aligned}
$$

Consider perturbations in the form of

$$
L_{1}=\left(\begin{array}{cc}
L_{1, X_{1}} & 0 \\
0 & 0
\end{array}\right), \text { where } L_{1, X_{1}}=\epsilon\left(\begin{array}{cc}
B & 0 \\
0 & \frac{B}{B}
\end{array}\right), \quad B=\left(\begin{array}{cccc}
0 & \ldots & 0 & 0 \\
\ldots & & & \\
0 & \ldots & 0 & 0 \\
0 & \ldots & 0 & 1
\end{array}\right)
$$

Clearly, $L_{1}$ is real in the sense $\left\langle L_{1} \bar{u}, \bar{v}\right\rangle=\overline{\left\langle L_{1} u, v\right\rangle}$. Let $L_{\#}=L+L_{1}$ and the decomposition $X=X_{1} \oplus X_{2}$ is still invariant under $J L_{\#}$ and orthogonal with respect to $L_{\#}$. Therefore, $\sigma\left(J_{+}\left(B_{+}+\epsilon B\right)\right) \subset \sigma\left(J L_{\#}\right)$. By direct computation, we obtain the matrix

$$
J_{+}\left(B_{+}+\epsilon B\right)=\left(\begin{array}{ccccc}
i \mu & 0 & \ldots & 0 & i \epsilon \mu{\overline{b_{1}}}^{-1} \\
1 & i \mu & \ldots & 0 & \epsilon{\overline{b_{1}}}^{-1} \\
\ldots & & & & \\
0 & 0 & \ldots & 1 & i \mu
\end{array}\right)
$$

and its characteristic polynomial

$$
\operatorname{det}\left(\lambda-J_{+}\left(B_{+}+\epsilon B\right)\right)=(-i)^{k} p(i(\lambda-i \mu))
$$

where

$$
p(\lambda)=\lambda^{k}-\epsilon b \lambda+\epsilon b \mu, \quad b=(-i)^{k-1}{\overline{b_{1}}}^{-1} \in\{ \pm 1\}
$$

To find hyperbolic eigenvalues of $J_{+}\left(B_{+}+\epsilon B\right)$, it is equivalent to show that $p(\lambda)=0$ has a root $\lambda \notin \mathbf{R}$. Choose the sign of $\epsilon$ such that $\epsilon b \mu>0$. Denote $c_{1}, \cdots, c_{k}$ to be
all the $k$-th roots of $-|b \mu|$, which are not real except for at most one. So we can assume $\operatorname{Im} c_{1} \neq 0$. Let $\delta=|\epsilon|^{\frac{1}{k}}$. Then $p(\lambda)=0$ is equivalent to

$$
\begin{equation*}
\left(\frac{\lambda}{\delta}\right)^{k}-\delta|b| \frac{|\mu|}{\mu}\left(\frac{\lambda}{\delta}\right)+|b \mu|=0 \tag{9.17}
\end{equation*}
$$

When $\delta \ll 1$, by the Implicit Function Theorem, (9.17) has $k$ roots of the form

$$
\frac{\lambda_{j}}{\delta}=c_{j}+O(\delta), j=1, \cdots, k
$$

among which $\lambda_{1}=\delta c_{1}+O\left(\delta^{2}\right)$ satisfies $\operatorname{Im} \lambda_{1} \neq 0$. This implies that $J_{+}\left(B_{+}+\epsilon B\right)$ has a hyperbolic eigenvalue of the form $i \mu-i \delta c_{1}+O\left(\delta^{2}\right)$.

Cases 3d: $\mu=0$ and $Y$ contains a non-trivial Jordan chain of length $\geq 3$. Let $u_{j}=(J L)^{j-1} u_{1}, j=1, \ldots, k,(k \geq 3)$ be a Jordan chain of $J L$ such that $\langle L \cdot, \cdot\rangle$ is non-degenerate on $\operatorname{span}\left\{u_{1}, \ldots, u_{k}\right\}$.

In this case one may proceed as in the above with

$$
p(\lambda)=\lambda^{k}-\epsilon b \lambda=\lambda\left(\lambda^{k-1}-\epsilon b\right)
$$

Choose $\epsilon$ such that $\epsilon b<0$. Since $k \geq 3, p(\lambda)$ has a complex root which implies that $J_{+}\left(B_{+}+\epsilon B\right)$ has a hyperbolic eigenvalue. For $\mu=0$ and $k=2$, by straightforward computations, it can be shown that $J$ must be degenerate and $J_{+}\left(B_{+}+\epsilon B\right)$ has only eigenvalue 0 and a purely imaginary eigenvalues for any $2 \times 2$ Hermitian matrix $B$ and $\epsilon \ll 1$.

Case 4: $i \mu \in \sigma(J L) \cap i \mathbf{R} \backslash\{0\}$ and $\langle L \cdot, \cdot\rangle$ is degenerate on $E_{i \mu} \neq\{0\}$.
In this case, Proposition 2.3 implies that $i \mu$ must be non-isolated in $\sigma(J L)$ and we start with the following lemma to isolate $i \mu$ through a perturbation.

Lemma 9.4. Assume (H1-3). Suppose $i \mu \in \sigma(J L) \cap i \mathbf{R}$ is non-isolated in $\sigma(J L)$. For any $\epsilon>0$, there exists a symmetric bounded linear operator $L_{1}: X \rightarrow$ $X^{*}$ satisfying (12.12) such that $\left|L_{1}\right|<\epsilon$ and $i \mu \in \sigma\left(J L_{\#}\right)$ is an isolated eigenvalue, where $L_{\#}=L+L_{1}$, and $\left\langle L_{\#} u, u\right\rangle>0$ for some generalized eigenvector $u$ of the eigenvalue $i \mu$ of $J L_{\#}$.

Proof. Since $\sigma(J L)$ is symmetric about both real and imaginary axes, without loss of generality we can assume that $\mu \geq 0$.

Let $X=\Sigma_{j=0}^{6} X_{j}$ be the decomposition given in Theorem 2.1 with associated projections $P_{j}$. We will use the notations there in the rest of the proof. Recall Theorem 2.1 is proved without the complexification, i.e. in the framework of real Hilbert space $X$ and real operators $J$ and $L$, the resulted decomposition and operators are real. After the complexification, $X_{j}$ are real in the sense of (9.15) and the operators satisfy (12.12) and the blocks in $L$ and $J$ are anti-linear.

As $i \mu \in \sigma(J L)$ is assumed to be non-isolated and $\operatorname{dim} X_{j}<\infty, j \neq 0,3$, it must hold $i \mu \in \sigma\left(A_{3}\right)$. Since $A_{3}$ is anti-self-adjoint with respect to the positive definite Hermitian form $\left\langle L_{X_{3}} \cdot, \cdot\right\rangle$, it induces a resolution of the identity. Namely there exists a family of projections $\left\{\Pi_{\lambda}\right\}_{\lambda \in \mathbf{R}}$ on $X_{3}$ such that
(1) $\lim _{\lambda \rightarrow \lambda_{0}+} \Pi_{\lambda} u=\Pi_{\lambda_{0}} u$, for all $\lambda_{0} \in \mathbf{R}$ and $u \in X_{3}$;
(2) $\Pi_{\lambda_{1}} \Pi_{\lambda_{2}}=\Pi_{\min \left\{\lambda_{1}, \lambda_{2}\right\}}$, for all $\lambda_{1,2} \in \mathbf{R}$;
(3) $\left\langle L_{X_{3}} \Pi_{\lambda} u_{1}, u_{2}\right\rangle=\left\langle L_{X_{3}} u_{1}, \Pi_{\lambda} u_{2}\right\rangle$ for any $u_{1,2} \in X_{3}$ and $\lambda \in \mathbf{R}$;
(4) $u=\int_{-\infty}^{+\infty} d \Pi_{\lambda} u, A_{3} u=\int_{-\infty}^{+\infty} i \lambda d \Pi_{\lambda} u$, for any $u \in X_{3}$;
(5) $d \Pi_{\lambda}=\overline{d \Pi_{-\lambda}}$ for any $\lambda \in \mathbf{R}$.

Here the last property is due to the fact that $J$ and $L$ are real satisfying (12.12). For $\mu>0$, define a perturbation $\tilde{L}_{1}: X_{3} \rightarrow X_{3}^{*}$ by

$$
\tilde{L}_{1} u=L_{X_{3}}\left(\int_{|\lambda-\mu|<\nu} \frac{\lambda-\mu}{\lambda} d \Pi_{\lambda} u+\int_{|\lambda+\mu|<\nu} \frac{\lambda+\mu}{\lambda} d \Pi_{\lambda} u\right), \quad \forall u \in X_{3}
$$

where $\nu \in(0, \mu)$ is a small constant to be determined later. As $\Pi_{\lambda} u$ is continuous from the right, the integrals take the same values on the open intervals or half open half closed interval like $[\mu-\nu, \mu+\nu)$.

For $\mu=0$, define

$$
\tilde{L}_{1} u=L_{X_{3}} \int_{-\nu}^{\nu} d \Pi_{\lambda} u, \quad \forall u \in X_{3}
$$

where again $\nu>0$ is determined later.
For $\mu>0$, like $L_{X_{3}}$, it is clear that $\tilde{L}_{1}$ is anti-linear satisfying (12.9). We will verify that $\tilde{L}_{1}$ is also real and symmetric. For any $u \in X$, one may compute using $d \Pi_{\lambda}=\overline{d \Pi_{-\lambda}}$

$$
\begin{aligned}
\overline{\tilde{L}_{1} u} & =L_{X_{3}}\left(\int_{|\lambda-\mu|<\nu} \frac{\lambda-\mu}{\lambda} d \Pi_{\lambda} u+\int_{|\lambda+\mu|<\nu} \frac{\lambda+\mu}{\lambda} d \Pi_{\lambda} u\right) \\
& =L_{X_{3}}\left(\int_{|\lambda-\mu|<\nu} \frac{\lambda-\mu}{\lambda} \overline{d \Pi_{\lambda}} \bar{u}+\int_{|\lambda+\mu|<\nu} \frac{\lambda+\mu}{\lambda} \overline{d \Pi_{\lambda}} \bar{u}\right) \\
& =L_{X_{3}}\left(\int_{|\lambda-\mu|<\nu} \frac{\lambda-\mu}{\lambda} d \Pi_{-\lambda} \bar{u}+\int_{|\lambda+\mu|<\nu} \frac{\lambda+\mu}{\lambda} d \Pi_{-\lambda} \bar{u}\right) .
\end{aligned}
$$

Through a change of variable $\lambda \rightarrow-\lambda$, we obtain $\bar{L}_{1} u=\tilde{L}_{1} \bar{u}$, namely, $\tilde{L}_{1}$ is real (the complxification of a real linear operator). Moreover, for any $u_{1,2} \in X$, we have

$$
\begin{aligned}
& \left\langle\tilde{L}_{1} u_{1}, u_{2}\right\rangle=\left\langle L_{X_{3}}\left(\int_{|\lambda-\mu|<\nu} \frac{\lambda-\mu}{\lambda} d \Pi_{\lambda} u_{1}+\int_{|\lambda+\mu|<\nu} \frac{\lambda+\mu}{\lambda} d \Pi_{\lambda} u_{1}\right), u_{2}\right\rangle \\
= & \int_{|\lambda-\mu|<\nu} \frac{\lambda-\mu}{\lambda} d\left\langle L_{X_{3}} \Pi_{\lambda} u_{1}, u_{2}\right\rangle+\int_{|\lambda+\mu|<\nu} \frac{\lambda+\mu}{\lambda} d\left\langle L_{X_{3}} \Pi_{\lambda} u_{1}, u_{2}\right\rangle .
\end{aligned}
$$

Since $L_{X_{3}}$ is Hermitian and $\left\langle\Pi_{\lambda} \cdot, \cdot\right\rangle=\left\langle\cdot, \Pi_{\lambda} \cdot\right\rangle$ on $X_{3}$, we obtain that $\tilde{L}_{1}$ is Hermitian. Therefore, $\left\langle\tilde{L}_{1} \cdot, \cdot\right\rangle$ is the complexification of a real bounded symmetric quadratic form on $X_{3}$. Clearly, in the equivalent norm $\left\langle L_{X_{3}} u, u\right\rangle^{\frac{1}{2}}$ on $X_{3}$,

$$
\left|\tilde{L}_{1}\right| \leq \frac{\nu}{\mu-\nu} \rightarrow 0, \quad \text { as } \nu \rightarrow 0
$$

The same properties also hold for $\tilde{L}_{1}$ for $\mu=0$ and we skip the details.
Let

$$
L_{1}=P_{3}^{*} \tilde{L}_{1} P_{3}, \quad L_{\#}=L-L_{1}
$$

Accordingly in this decomposition

$$
L_{\#} \longleftrightarrow\left(\begin{array}{ccccccc}
0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & B_{14} & 0 & 0 \\
0 & 0 & L_{X_{2}} & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & L_{X_{3}}+\tilde{L}_{1} & 0 & 0 & 0 \\
0 & B_{14}^{*} & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & B_{56} \\
0 & 0 & 0 & 0 & 0 & B_{56}^{*} & 0
\end{array}\right)
$$

From Corollary 2.1, one can compute

$$
J L_{\#} \longleftrightarrow\left(\begin{array}{ccccccc}
0 & A_{01} & A_{02} & A_{03}\left(I-L_{X_{3}}^{-1} \tilde{L}_{1}\right) & A_{04} & 0 & 0 \\
0 & A_{1} & A_{12} & A_{13}\left(I-L_{X_{3}}^{-1} \tilde{L}_{1}\right) & A_{14} & 0 & 0 \\
0 & 0 & A_{2} & 0 & A_{24} & 0 & 0 \\
0 & 0 & 0 & A_{3}\left(I-L_{X_{3}}^{-1} \tilde{L}_{1}\right) & A_{34} & 0 & 0 \\
0 & 0 & 0 & 0 & A_{4} & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & A_{5} & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & A_{6}
\end{array}\right) .
$$

Due to the upper triangular structure of $J L_{\#}$ and the finite dimensionality of $X_{1,2}$, in order to prove that $i \mu$ belongs to and is isolated in $\sigma\left(J L_{\#}\right)$, it suffices to show that $i \mu$ belongs to and is isolated in $\sigma\left(A_{3}\left(I-L_{X_{3}}^{-1} \tilde{L}_{1}\right)\right)$. In fact, for any $u \in X_{3}$,

$$
\begin{equation*}
A_{3}\left(I-L_{X_{3}}^{-1} \tilde{L}_{1}\right) u=i \int_{S} \mu d \Pi_{\lambda} u+i \int_{\mathbf{R} \backslash S} \lambda d \Pi_{\lambda} u \tag{9.18}
\end{equation*}
$$

where $S=(-\mu-\nu,-\mu+\nu) \cup(\mu-\nu, \mu+\nu)$. Since $\Pi_{\lambda}$ is not constant on $S$ as $i \mu \in \sigma\left(A_{3}\right)$, we obtain that $i \mu$ is an isolated eigenvalue of $A_{3}\left(I-L_{X_{3}}^{-1} \tilde{L}_{1}\right)$ and thus of $\sigma\left(J L_{\#}\right)$ as well. Indeed, for any $u \in R\left(\Pi_{\mu+\nu}-\Pi_{\mu-\nu}\right)$, by (9.18) we have

$$
A_{3}\left(I-L_{X_{3}}^{-1} \tilde{L}_{1}\right) u=i \mu u
$$

So $i \mu$ is an eigenvalues of $A_{3}\left(I-L_{X_{3}}^{-1} \tilde{L}_{1}\right)$. To show $i \mu$ is isolated, taking any $\alpha \in \mathbf{C}$ such that $0<|\alpha-i \mu|<\nu$, then we have

$$
\left(\alpha-A_{3}\left(I-L_{X_{3}}^{-1} \tilde{L}_{1}\right)\right)^{-1}=\int_{S}(\alpha-i \mu)^{-1} d \Pi_{\lambda}+\int_{\mathbf{R} \backslash S}(\alpha-i \lambda)^{-1} d \Pi_{\lambda}
$$

which is clearly a bounded operator.
Finally, we prove that there exists a generalized eigenvector $u$ of $i \mu$ of $J L_{\#}$ such that $\left\langle L_{\#} u, u\right\rangle>0$. Since $\operatorname{dim} X_{1}<\infty$, there exists an integer $K>0$ such that

$$
X_{1}=Y_{\mu} \oplus \tilde{Y}, \quad \text { where } Y_{\mu}=\operatorname{ker}\left(A_{1}-i \mu\right)^{K} \cap X_{1}, \tilde{Y}=\left(A_{1}-i \mu\right)^{K} X_{1}
$$

In the following we proceed in the case of $\mu>0$ first. Let

$$
Z_{\mu}=\left\{u-(-i \mu)^{-K} P_{0}(J L-i \mu)^{K} u \mid u \in Y_{\mu}\right\}, \quad \tilde{Z}=\operatorname{ker} L \oplus \tilde{Y}
$$

Note that the upper triangular structure of $J L$ implies that

$$
\left(A_{1}-i \mu\right)^{K}=\left.P_{1}(J L-i \mu)^{K}\right|_{X_{1}} .
$$

Using this observation and the invariance of $\tilde{X}=\operatorname{ker} L \oplus X_{1}$ under $J L$, we obtain through straightforward computations

$$
\begin{equation*}
\tilde{X}=Z_{\mu} \oplus \tilde{Z}, Z_{\mu}=\operatorname{ker}(J L-i \mu)^{K} \cap \tilde{X}, \tilde{Z}=(J L-i \mu)^{K} \tilde{X} \tag{9.19}
\end{equation*}
$$

and on the invariant subspaces $Z_{\mu}$ and $\tilde{Z}$

$$
\begin{equation*}
\sigma\left(\left.J L\right|_{Z_{\mu}}\right)=\{i \mu\} \quad \text { if } \quad i \mu \in \sigma\left(A_{1}\right), \quad i \mu \notin \sigma\left(\left.J L\right|_{\tilde{Z}}\right) \tag{9.20}
\end{equation*}
$$

Let $P: \tilde{X} \rightarrow \tilde{Z}$ be the projection associated to the above decomposition and $u_{3} \in X_{3}$ be such that

$$
A_{3}\left(I-L_{X_{3}}^{-1} \tilde{L}_{1}\right) u_{3}=i \mu u_{3}
$$

The structure of $J L_{\#}$ implies $\left(J L_{\#}-i \mu\right) u_{3} \in \tilde{X}$. Let

$$
\tilde{u}=\left(\left.(J L-i \mu)\right|_{\tilde{Z}}\right)^{-1} \tilde{P}\left(J L_{\#}-i \mu\right) u_{3} \in \tilde{Z} \subset \tilde{X}, \quad u=u_{3}-\tilde{u}
$$

By using $\left.\left(L_{\#}-L\right)\right|_{\tilde{X}}=0$, it is easy to verify that

$$
\left(J L_{\#}-i \mu\right) u \in Z_{\mu}
$$

which implies

$$
\left(J L_{\#}-i \mu\right)^{K+1} u=0
$$

From the structure of $L_{\#}$, straightforward computation leads to

$$
\left\langle L_{\#} u, u\right\rangle=\left\langle\left(L_{3}+\tilde{L}_{1}\right) u_{3}, u_{3}\right\rangle>0
$$

for $0<\nu \ll 1$.
The case of $\mu=0$ is largely similar. Let

$$
Z_{0}=Y_{0} \oplus \operatorname{ker} L, \quad \tilde{Z}=\tilde{Y}
$$

and (9.19) and (9.20) still hold. The rest of the argument follows in exactly the same procedure.

We return to construct a perturbation $L_{1}$ to $L$ to create unstable eigenvalues. In Case $4, E^{D}$ in Proposition 2.2 is non-trivial and finite dimensional, therefore

$$
\begin{equation*}
\exists 0 \neq u_{0} \in \operatorname{ker}(J L-i \mu) \text { such that }\left\langle L u_{0}, u_{0}\right\rangle=0 \tag{9.21}
\end{equation*}
$$

where $u_{0} \in E^{D}$. Since $\mu \neq 0$ implies $\overline{u_{0}} \in E_{-i \mu}$ with $\left\langle L u_{0}, \overline{u_{0}}\right\rangle=0$, let

$$
\begin{equation*}
Y_{0}=\operatorname{span}\left\{u_{0}, \overline{u_{0}}\right\} \subset \operatorname{ker}(J L-i \mu) \oplus \operatorname{ker}(J L+i \mu) \tag{9.22}
\end{equation*}
$$

The following decomposition lemma is our first step in the construction of a hyperbolically generating perturbation.

Lemma 9.5. Suppose $0 \neq i \mu \in \sigma(J L) \cap i \mathbf{R}$ satisfying (9.21). Let $Y_{0}$ be defined in (9.22). Then there exists $w \in D(J L)$ with $\bar{w} \neq w$ and a codim-4 closed subspace $Y_{1} \subset X$ satisfying (9.15) such that $X=Y_{0} \oplus Y_{1} \oplus Y_{2}$, where $Y_{2}=\operatorname{span}\{w, \bar{w}\}$. Moreover, in this decomposition and the bases $\left\{u_{0}, \overline{u_{0}}\right\},\{w, \bar{w}\}$ on $Y_{0,2}$ respectively, $L$ and $J L$ take the forms

$$
L \longleftrightarrow\left(\begin{array}{ccc}
0 & 0 & I_{2 \times 2} \\
0 & L_{Y_{1}} & 0 \\
I_{2 \times 2} & 0 & 0
\end{array}\right), J L \longleftrightarrow\left(\begin{array}{ccc}
i \mu \Lambda & A_{01} & A_{02} \\
0 & A_{1} & A_{12} \\
0 & 0 & i \mu \Lambda
\end{array}\right), \Lambda=\left(\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right)
$$

Here, all blocks are bounded operators except $A_{1}=J_{Y_{1}} L_{Y_{1}}$ and $\left(Y_{1}, J_{Y_{1}}, L_{Y_{1}}\right)$ satisfies (H1-3).

Proof. Let

$$
\tilde{Y}=\left\{u \in X \mid\left\langle L u, u_{0}\right\rangle=0=\left\langle L u, \overline{u_{0}}\right\rangle\right\} \supset\left\{u_{0}, \overline{u_{0}}\right\}
$$

Clearly, $\tilde{Y}$, satisfying (9.15), is the complexification of some real codim-2 subspace. Lemma 3.2 implies that $\tilde{Y}$ is invariant under $J L$. Let

$$
\tilde{Y}_{1}=\left\{u \in \tilde{Y} \mid\left(u, u_{0}\right)=\left(u, \overline{u_{0}}\right)=0\right\}
$$

Since $Y_{0} \cap \operatorname{ker} L=\{0\}$ and $D(J L)$ is dense in $X$, there exists a 2-dim subspace $\tilde{Y}_{2} \subset D(J L)$ such that $\langle L u, v\rangle, u \in Y_{0}$ and $v \in \tilde{Y}_{2}$, defines a non-degenerate bilinear form on $Y_{0} \otimes \tilde{Y}_{2}$. Clearly, we have $X=Y_{0} \oplus \tilde{Y}_{1} \oplus \tilde{Y}_{2}$ and in this decomposition $L$ takes the form

$$
L \longleftrightarrow\left(\begin{array}{ccc}
0 & 0 & B_{02} \\
0 & L_{Y_{1}} & B_{12} \\
B_{02}^{*} & B_{12}^{*} & B_{22}
\end{array}\right)
$$

where $B_{02}: \tilde{Y}_{2} \rightarrow Y_{0}^{*}$ is non-degenerate and $B_{22}^{*}=B_{22}$. Through exactly the same procedure as in the proof of Proposition 6.1, we may obtain subspaces $Y_{1}$ and $\tilde{X}_{2}$ as graphs of bounded linear operators from $\tilde{Y}_{1,2}$ to $Y_{0}$ such that $X=Y_{0} \oplus Y_{1} \oplus X_{2}$ and in this decomposition $L$ takes the form

$$
L \longleftrightarrow\left(\begin{array}{ccc}
0 & 0 & B \\
0 & L_{Y_{1}} & 0 \\
B^{*} & 0 & 0
\end{array}\right)
$$

where $B: X_{2} \rightarrow Y_{0}^{*}$ is non-degenerate. There exists $w \in X_{2} \subset D(J L)$ such that $\left\langle L u_{0}, w\right\rangle=1$ and $\left\langle L u_{0}, \bar{w}\right\rangle=\left\langle L \overline{u_{0}}, w\right\rangle=0$, which also implies $\left\langle L \overline{u_{0}}, \bar{w}\right\rangle=0$, where (12.12) is used. Let $Y_{2}=\operatorname{span}\{w, \bar{w}\}$. From the definition of $w, \tilde{Y}$, and $Y_{1}$, we have $X=Y_{0} \oplus Y_{1} \oplus Y_{2}$, associated with projections $P_{0,1,2}$, and in this decomposition, the desired block form of $L$ is achieved. Applying Lemma 12.3 to $X=\left(Y_{0} \oplus Y_{2}\right) \oplus Y_{1}$, we obtain that $\left(Y_{1}, J_{Y_{1}}, L_{Y_{1}}\right)$ satisfies (H1-3), where $J_{Y_{1}}=P_{1} J P_{1}^{*}$. The upper triangular block form of $J L$ is due to the invariance of $Y_{0}$ and $\tilde{Y}=Y_{0} \oplus Y_{2}$.

To complete the proof of the lemma, we are left to show $P_{2} J L w=i \mu w$, which along with the facts that $\tilde{Y}$ satisfies (9.15) and $J L$ satisfies (12.12) also implies $P_{2} J L \bar{w}=-i \mu \bar{w}$. From $(J L)^{*}=-L J$ (Corollary 12.1), we have

$$
\left\langle L J L w, u_{0}\right\rangle=-\left\langle L w, J L u_{0}\right\rangle=i \mu\left\langle L w, u_{0}\right\rangle=i \mu
$$

and similarly $\left\langle L J L w, \overline{u_{0}}\right\rangle=-i \mu\left\langle L w, \overline{u_{0}}\right\rangle=0$. According to the definitions of $\tilde{Y}$ and $w$, we obtain $P_{2} J L w=i \mu w$ and the lemma is proved.

With the above lemmas, we are ready to construct a perturbed energy operator $L_{\#}$ to create unstable eigenvalues of $J L_{\#}$ near $i \mu$ in the Case 4 . We start with the decomposition given in Lemma 9.5. Since $i \mu$ is an eigenvalue of $J L$ non-isolated in $\sigma(J L)$, we have $i \mu \in \sigma\left(J_{Y_{1}} L_{Y_{1}}\right)$ and is non-isolated in $\sigma\left(J_{Y_{1}} L_{Y_{1}}\right)$. From Lemma 9.4, there exists a sufficiently small symmetric bounded linear operator $\tilde{L}_{2}: Y_{1} \rightarrow Y_{1}^{*}$ such that $i \mu \in \sigma\left(J_{Y_{1}}\left(L_{Y_{1}}+\tilde{L}_{2}\right)\right)$ and is isolated with an eigenvector $u_{1} \in Y_{1}$ satisfying $\left\langle\left(L_{Y_{1}}+\tilde{L}_{2}\right) u_{1}, u_{1}\right\rangle>0$. Let $L_{2}=P_{1}^{*} \tilde{L}_{2} P_{1}$ and $\tilde{L}_{\#}=L+L_{2}$, then the block forms of $L_{\#}$ and $J \tilde{L}_{\#}$ imply that $i \mu \in \sigma\left(J \tilde{L}_{\#}\right)$ is isolated and

$$
\begin{equation*}
u_{0}, u_{1} \in \operatorname{ker}\left(J \tilde{L}_{\#}-i \mu\right), \quad\left\langle\tilde{L}_{\#} u_{0}, u_{0}\right\rangle=0, \quad \text { and }\left\langle\tilde{L}_{\#} u_{1}, u_{1}\right\rangle>0 \tag{9.23}
\end{equation*}
$$

Since $i \mu$ is isolated in $\sigma\left(J \tilde{L}_{\#}\right)$, Proposition 2.3 implies that $\left\langle\tilde{L}_{\#} \cdot, \cdot\right\rangle$ is non-degenerate on $E_{i \mu}\left(J \tilde{L}_{\#}\right)$, the subspace of generalized eigenvectors of $i \mu$ for $J \tilde{L}_{\#}$. Moreover,
by $(9.23),\left\langle\tilde{L}_{\#} \cdot, \cdot\right\rangle$ is sign indefinite on $E_{i \mu}\left(J \tilde{L}_{\#}\right)$. This situation has been covered in Case 3. Therefore, there exists a sufficient small symmetric bounded linear operator $L_{3}: X \rightarrow X^{*}$ such that there exists $\lambda \in \sigma\left(J L_{\#}\right) \backslash i \mathbf{R}$ sufficiently close to $i \mu$, where $L_{\#}=L+L_{2}+L_{3}$.

Case 5: $i \mu \in \sigma(J L) \cap i \mathbf{R} \backslash\{0\}$ is non-isolated and $\langle L \cdot, \cdot\rangle$ is negative definite on $E_{i \mu} \neq\{0\}$.

Much as in Case 4 (but more easily), we can construct sufficiently small symmetric bounded perturbations to the energy operator $L$ to create unstable eigenvalues. In fact, Proposition 2.2 implies that in Case 5, it holds $\operatorname{ker}(J L-i \mu)=E_{i \mu}$. Let

$$
Y_{0}=E_{i \mu} \oplus E_{-i \mu}, \quad Y=Y_{0}^{\perp_{L}}=\left\{v \in X \mid\langle L v, u\rangle=\langle L v, \bar{u}\rangle=0, \forall u \in E_{i \mu}\right\}
$$

Since $\langle L \cdot, \cdot\rangle$ is negative on $Y_{0}$, Lemma 12.2 implies that $X=Y_{0} \oplus Y$ associated with projections $P_{Y_{0}, Y}$. In this decomposition $L$ and $J L$ take the forms

$$
L \longleftrightarrow\left(\begin{array}{cc}
L_{Y_{0}} & 0 \\
0 & L_{Y}
\end{array}\right), \quad J L \longleftrightarrow\left(\begin{array}{cc}
A_{0} & 0 \\
0 & A
\end{array}\right)
$$

where $A_{0}$ is a bounded operator satisfying $A^{2}+\mu^{2}=0$. Lemma 12.3 implies that $A=J_{Y} L_{Y}$ and $\left(Y, J_{Y}, L_{Y}\right)$ satisfies (H1-3). Clearly, it still holds that $i \mu \in \sigma\left(J_{Y} L_{Y}\right)$ and is non-isolated there. Applying Lemma 9.4 to $L_{Y}$, we obtain a perturbation $\tilde{L}: Y \rightarrow Y^{*}$ such that $i \mu$ is an isolated point in $\sigma\left(J_{Y}\left(L_{Y}+\tilde{L}\right)\right)$. Let $\tilde{L}_{\#}=L+P_{Y}^{*} \tilde{L} P_{Y}$ and we obtain that $i \mu$ is an isolated point in $\sigma\left(J \tilde{L}_{\#}\right)$ with $\left\langle\tilde{L}_{\#} \cdot, \cdot\right\rangle$ sign indefinite on its eigenspace. This is a case covered in Case 3 and thus there exists a sufficient small symmetric bounded linear perturbation $L_{\#}$ to $L$ so that $J L_{\#}$ has an unstable eigenvalue close to $i \mu$.

Proof of Theorem 2.6. It suffices to show that Cases 3, 4, 5 cover all the cases in Theorem 2.6. In fact, if $\langle L \cdot, \cdot\rangle$ is degenerate on $E_{i \mu} \neq\{0\}$ and $\mu \neq 0$, this is precisely Case 4. Let us consider the case when $\langle L \cdot, \cdot\rangle$ is non-degenerate on $E_{i \mu}$ and satisfies the assumptions in Theorem 2.6. Then $\langle L \cdot, \cdot\rangle$ is either sign indefinite on $E_{i \mu}$ (Case 3) or is negative definite on $E_{i \mu}$ for an eigenvalue $i \mu \neq 0$ non-isolated in $\sigma(J L)$ (Case 5).

## CHAPTER 10

## Proof of Theorem 2.7 where (H2.b) is weakened

In this chapter, we consider the case when (H2.b) is weakened, namely, $L$ is only assumed to be positive on $X_{+}$, but not necessarily uniformly positive. More precisely, we will prove Theorem 2.7 under hypotheses (B1-5) given in Section 2.6. In Section 11.6, as an example we will consider the stability of traveling waves of a nonlinear Schrödinger equation with non-vanishing condition at infinity in two dimensions.

Initial decomposition of the phase space. We adopt the notations as in Chapter 3. Let $P_{ \pm, 0}: X \rightarrow X_{ \pm, 0}$ be the projections associated to the decomposition $X=X_{-} \oplus \operatorname{ker} L \oplus X_{+}$, where $X_{0}=\operatorname{ker} L$, and

$$
\tilde{X}_{ \pm, 0}^{*}=P_{ \pm, 0}^{*} X_{ \pm, 0}^{*} \subset X^{*}
$$

We also let

$$
X_{\leq 0}=X_{-} \oplus \operatorname{ker} L, \quad P_{\leq 0}=P_{0}+P_{-}=I-P_{+}, \quad \tilde{X}_{\leq 0}^{*}=\tilde{X}_{-}^{*} \oplus \tilde{X}_{0}^{*}
$$

Clearly, we have

$$
\begin{equation*}
\tilde{X}_{+}^{*}=\operatorname{ker} i_{X_{\leq 0}}^{*}, \quad \tilde{X}_{\leq 0}^{*}=\operatorname{ker} i_{X_{+}}^{*} \subset Q_{0}(X), \quad X^{*}=\tilde{X}_{\leq 0}^{*} \oplus \tilde{X}_{+}^{*} \tag{10.1}
\end{equation*}
$$

where assumption (B5) is used. Since $\langle L u, u\rangle<0$ on $X_{-} \backslash\{0\}$ and $\operatorname{dim} X_{-}=$ $n^{-}(L)<\infty$, there exists $\delta>0$ such that

$$
\langle L u, u\rangle \leq-\delta\|u\|^{2}, \quad \forall u \in X_{-} .
$$

From (B4), we also have

$$
L X_{+} \subset \tilde{X}_{+}^{*}, \quad L X_{\leq 0}=\tilde{X}_{-}^{*} \subset \tilde{X}_{\leq 0}^{*}
$$

Denote

$$
L_{+}=i_{X_{+}}^{*} L i_{X_{+}}: X_{+} \rightarrow X_{+}^{*}, \quad L_{\leq 0}=i_{X_{\leq 0}}^{*} L i_{X_{\leq 0}}: X_{\leq 0} \rightarrow X_{\leq 0}^{*}
$$

which along with the $L$-orthogonality in ( $\mathbf{B 4}$ ) implies

$$
L=P_{+}^{*} L_{+} P_{+}+P_{\leq 0}^{*} L_{\leq 0} P_{\leq 0}
$$

While the decomposition is not necessarily $Q_{0}$-orthogonal, we have the following lemma. Let

$$
\begin{aligned}
& Q_{0}^{\leq 0,+}=i_{\leq 0}^{*} Q_{0} i_{X_{+}}: X_{+} \rightarrow X_{\leq 0}^{*}, \quad Q_{0}^{+, \leq 0}=i_{X_{+}}^{*} Q_{0} i_{\leq 0}: X_{\leq 0} \rightarrow X_{+}^{*} \\
& Q_{0}^{\leq 0}=i_{X_{\leq 0}}^{*} Q_{0} i_{X_{\leq 0}}: X_{\leq 0} \rightarrow X_{\leq 0}^{*}, \quad Q_{0}^{+}=i_{X_{+}}^{*} Q_{0} i_{X_{+}}: X_{+} \rightarrow X_{+}^{*}
\end{aligned}
$$

Clearly, $Q_{0}^{\leq 0,+}=\left(Q_{0}^{+, \leq 0}\right)^{*}$ and in the decomposition $X=X_{\leq 0} \oplus X_{+}$and $X^{*}=$ $P_{\leq 0}^{*} X_{\leq 0}^{*} \oplus P_{+}^{*} X_{+}^{*}$, operator $Q_{0}$ takes the form $\left(\begin{array}{cc}Q_{0}^{\leq 0} & Q_{0}^{\leq 0,+} \\ Q_{0}^{+, \leq 0} & Q_{0}^{+}\end{array}\right)$. Since $\left\langle Q_{0} u, u\right\rangle>0$ for all $0 \neq u \in X, Q_{0}^{+}$and $Q_{0}^{\leq 0}$, as well as $L_{+}$, are bounded, symmetric, and
positive. Therefore, $Q_{0}^{\leq 0}: X_{\leq 0} \rightarrow X_{\leq 0}^{*}$ and $Q_{0}^{+}, L_{+}: X_{+} \rightarrow X_{+}^{*}$ are injective with dense ranges. Consequently, $\left(Q_{0}^{\leq 0}\right)^{-1}: X_{\leq 0}^{*} \rightarrow X_{\leq 0}$ and $\left(Q_{0}^{+}\right)^{-1}, L_{+}^{-1}: X_{+}^{*} \rightarrow X_{+}$ are densely defined, closed, and positive operators with

$$
\left(\left(Q_{0}^{\leq 0}\right)^{-1}\right)^{*}=\left(Q_{0}^{\leq 0}\right)^{-1},\left(\left(Q_{0}^{+}\right)^{-1}\right)^{*}=\left(Q_{0}^{+}\right)^{-1}
$$

and $\left(L_{+}^{-1}\right)^{*}=L_{+}^{-1}$.
Lemma 10.1. It holds that $P_{+}^{*} Q_{0}^{+}\left(X_{+}\right) \subset \tilde{X}_{+}^{*}$ is dense in $\tilde{X}_{+}^{*}$ and

$$
Q_{0}(X)=\tilde{X}_{\leq 0}^{*} \oplus P_{+}^{*} Q_{0}^{+}\left(X_{+}\right), \quad P_{+}^{*} Q_{0}^{+}\left(X_{+}\right)=Q_{0}(X) \cap \tilde{X}_{+}^{*}
$$

with $\left(Q_{0}^{\leq 0}\right)^{-1}$ and $\left(Q_{0}^{+}\right)^{-1} Q_{0}^{+, \leq 0}$ being bounded operators. Moreover,

$$
A \triangleq Q_{0}^{-1} P_{+}^{*} Q_{0}^{+}: X_{+} \rightarrow X_{2}, \quad \text { where } X_{2}=Q_{0}^{-1}\left(\tilde{X}_{+}^{*}\right) \subset X
$$

is an isomorphism.
This lemma makes the natural connection between $Q_{0}(X)$ and $Q_{0}^{+}\left(X_{+}\right)$.
Proof. Since the quadratic form $\left\langle Q_{0} u, u\right\rangle$ is positive on $X$, we have that $Q_{0}$ : $X \rightarrow X^{*}$ is injective with dense $Q_{0}(X) \subset X^{*}$. As $\tilde{X}_{\leq 0}^{*}=\operatorname{ker} i_{X_{+}}^{*} \subset Q_{0}(X)$ due to (B5) and $X^{*}=\tilde{X}_{\leq 0}^{*} \oplus \tilde{X}_{+}^{*}$, we obtain that $\tilde{X}_{+}^{*} \cap Q_{0}(X)$ is dense in $\tilde{X}_{+}^{*}$ and $Q_{0}(X)=\tilde{X}_{\leq 0}^{*} \oplus\left(Q_{0}(\bar{X}) \cap \tilde{X}_{+}^{*}\right)$. In the rest of the proof, we study $Q_{0}(X) \cap \tilde{X}_{+}^{*}$ and its associated properties.

Let $X_{1}=Q_{0}^{-1}\left(\tilde{X}_{\leq 0}^{*}\right) \subset X$, which is a closed subspace. Since $\tilde{X}_{\leq 0}^{*} \subset Q_{0}(X)$ and $Q_{0}$ is injective, $Q_{0}: X_{1} \rightarrow \tilde{X}_{\leq 0}^{*}$ is bounded, injective, and surjective and thus an isomorphism. Let

$$
\phi=\left(\left.Q_{0}\right|_{X_{1}}\right)^{-1} P_{\leq 0}^{*}: X_{\leq 0}^{*} \rightarrow X_{1}, \quad \phi_{\leq 0}=P_{\leq 0} \phi, \quad \phi_{+}=P_{+} \phi
$$

which are bounded operators. For any $f, g \in X_{\leq 0}^{*}$, since

$$
\left\langle g, \phi_{\leq 0} f\right\rangle=\left\langle P_{\leq 0}^{*} g, \phi f\right\rangle=\left\langle Q_{0} \phi g, \phi f\right\rangle,
$$

we obtain that $\phi_{\leq 0}: X_{\leq 0}^{*} \rightarrow X_{\leq 0}$ is symmetric and $\left\langle f, \phi_{\leq 0} f\right\rangle>0$ for any $0 \neq f \in$ $X_{\leq 0}^{*}$. Therefore $\phi_{\leq 0}^{-1}$ is a densely defined closed operator satisfying $\left(\phi_{\leq 0}^{-1}\right)^{*}=\phi_{\leq 0}^{-1}>$ 0 .

For any $f \in X_{\leq 0}^{*}$, let

$$
\phi f=u_{\leq 0}+u_{+}, \quad u_{+}=\phi_{+} f, \quad u_{\leq 0}=\phi_{\leq 0} f
$$

then we have

$$
Q_{0}^{\leq 0} u_{\leq 0}+Q_{0}^{\leq 0,+} u_{+}=f, \quad Q_{0}^{+, \leq 0} u_{\leq 0}+Q_{0}^{+} u_{+}=0
$$

It implies that $Q_{0}^{+, \leq 0} u_{\leq 0} \in Q_{0}^{+}\left(X_{+}\right)$and $u_{+}=-\left(Q_{0}^{+}\right)^{-1} Q_{0}^{+, \leq 0} u_{\leq 0}$. Therefore,

$$
\left(Q_{0}^{\leq 0}-Q_{0}^{\leq 0,+}\left(Q_{0}^{+}\right)^{-1} Q_{0}^{+, \leq 0}\right) u_{\leq 0}=f
$$

which implies that the closed positive symmetric operator $\phi_{\leq 0}^{-1}$ satisfies

$$
0<\phi_{\leq 0}^{-1}=Q_{0}^{\leq 0}-Q_{0}^{\leq 0,+}\left(Q_{0}^{+}\right)^{-1} Q_{0}^{+, \leq 0} \leq Q_{0}^{\leq 0}
$$

Here we also used $Q_{0}^{\leq 0,+}=\left(Q_{0}^{+, \leq 0}\right)^{*}$ and the positivity of the symmetric closed operator $\left(Q_{0}^{+}\right)^{-1}$. Therefore, $\phi_{\leq 0}$ is an isomorphism and

$$
\left(Q_{0}^{+}\right)^{-1} Q_{0}^{+, \leq 0}=-\phi_{+} \phi_{\leq 0}^{-1}
$$

is bounded. The above inequality also implies the boundedness of $\left(Q_{0}^{\leq 0}\right)^{-1} \leq \phi_{\leq 0}$.

On the one hand, for any $u \in X_{+}$, using $I=i_{X_{\leq 0}} P_{\leq 0}+i_{X_{+}} P_{+}$we can write

$$
P_{+}^{*} Q_{0}^{+} u=Q_{0} u-P_{\leq 0}^{*} i_{X_{\leq 0}}^{*} Q_{0} u=Q_{0}\left(I-\phi i_{X_{\leq 0}}^{*} Q_{0}\right) u
$$

Therefore, $P_{+}^{*} Q_{0}^{+}\left(X_{+}\right) \subset Q_{0}(X) \cap \tilde{X}_{+}^{*}$ and

$$
A \triangleq Q_{0}^{-1} P_{+}^{*} Q_{0}^{+}=I-\phi i_{X_{\leq 0}}^{*} Q_{0}: X_{+} \rightarrow X_{2}
$$

is bounded, where $X_{2}=Q_{0}^{-1}\left(\tilde{X}_{+}^{*}\right)$ is a closed subspace of $X$ and $Q_{0}(X) \cap \tilde{X}_{+}^{*}=$ $Q_{0}\left(X_{2}\right)$.

On the other hand, suppose $u=u_{\leq 0}+u_{+} \in X_{2}$, let $f=i_{X_{+}}^{*} Q_{0} u \in X_{+}^{*}$ and $f_{+}=P_{+}^{*} f=Q_{0} u \in \tilde{X}_{+}^{*}$. We have

$$
Q_{0}^{\leq 0} u_{\leq 0}+Q_{0}^{\leq 0,+} u_{+}=0, \quad Q_{0}^{+, \leq 0} u_{\leq 0}+Q_{0}^{+} u_{+}=f,
$$

and thus $u_{\leq 0}=-\left(Q_{0}^{\leq 0}\right)^{-1} Q_{0}^{\leq 0,+} u_{+}$. Substituting it into the second equation in the above, we obtain

$$
f=\left(Q_{0}^{+}-Q_{0}^{+, \leq 0}\left(Q_{0}^{\leq 0}\right)^{-1} Q_{0}^{\leq 0,+}\right) u_{+}=Q_{0}^{+} \tilde{u}_{+}
$$

where, from the above boundedness of $\left(Q_{0}^{+}\right)^{-1} Q_{0}^{+, \leq 0}$,

$$
\tilde{u}_{+}=\left(I-\left(Q_{0}^{+}\right)^{-1} Q_{0}^{+, \leq 0}\left(Q_{0}^{\leq 0}\right)^{-1} Q_{0}^{\leq 0,+}\right) u_{+} \in X_{+}
$$

It implies $f \in Q_{0}^{+}\left(X_{+}\right)$and thus $f_{+} \in P_{+}^{*} Q_{0}^{+}\left(X_{+}\right)$. Therefore

$$
Q_{0}(X) \cap \tilde{X}_{+}^{*} \subset P_{+}^{*} Q_{0}^{+}\left(X_{+}\right)
$$

Moreover, the above equality on $\tilde{u}_{+}$also implies

$$
A\left(I-\left(Q_{0}^{+}\right)^{-1} Q_{0}^{+, \leq 0}\left(Q_{0}^{\leq 0}\right)^{-1} Q_{0}^{\leq 0,+}\right) P_{+} u=Q_{0}^{-1} P_{+}^{*} Q_{0}^{+} \tilde{u}_{+}=Q_{0}^{-1} f_{+}=u
$$

Therefore we obtain

$$
A^{-1}=\left(I-\left(Q_{0}^{+}\right)^{-1} Q_{0}^{+, \leq 0}\left(Q_{0}^{\leq 0}\right)^{-1} Q_{0}^{\leq 0,+}\right) P_{+}
$$

is bounded and the proof of the lemma is complete.
Construction of $Y$. As our main concern is that $L_{+}$is not uniformly positive definite on $X_{+}$, we will actually work on the completion $Y_{+}$of $X_{+}$under the positive quadratic form $\left\langle L_{+}, \cdot\right\rangle$.

We start with a resolution of identity to rewrite $L_{+}$on $X_{+}$. From (B3), there exists $a>0$ such that

$$
\begin{equation*}
\frac{1}{C}\|u\|^{2} \leq\|u\|_{L_{+}, a}^{2} \leq C\|u\|^{2}, \quad \forall u \in X_{+} \tag{10.2}
\end{equation*}
$$

for some $C>0$, where, for $u, v \in X_{+},\|u\|_{L_{+}, a}^{2}=(u, u)_{L_{+}, a}$ and

$$
(u, v)_{L_{+}, a} \triangleq\left\langle\left(L_{+}+a Q_{0}^{+}\right) u, v\right\rangle=\left\langle\left(L+a Q_{0}\right) u, v\right\rangle .
$$

For $u, v \in X_{+}$, let

$$
\mathbb{L}=\left(L_{+}+a Q_{0}^{+}\right)^{-1} L_{+}: X_{+} \rightarrow X_{+}
$$

which implies $(\mathbb{L} u, v)_{L_{+}, a}=\langle L u, v\rangle$ and

$$
\begin{equation*}
\mathbb{D}=\left(Q_{0}^{+}\right)^{-1}\left(L_{+}+a Q_{0}^{+}\right): X_{+} \supset D(\mathbb{D})=\left(L_{+}+a Q_{0}^{+}\right)^{-1} Q_{0}^{+}\left(X_{+}\right) \rightarrow X_{+} \tag{10.3}
\end{equation*}
$$

Clearly, the Riesz representation $\mathbb{L}$ of $L_{+}$with respect to the equivalent metric $(\cdot, \cdot)_{L_{+}, a}$ is a bounded symmetric linear operator. Since

$$
\begin{equation*}
\mathbb{D}^{-1}=\left(L_{+}+a Q_{0}^{+}\right)^{-1} Q_{0}^{+}=a^{-1}(I-\mathbb{L}) \tag{10.4}
\end{equation*}
$$

is a bounded linear operator symmetric (and positive) with respect to $(\cdot, \cdot)_{L_{+}, a}, \mathbb{D}$ is self-adjoint with respect to $(\cdot, \cdot)_{L_{+}, a}$. In applications, if $Q_{1}$ is a uniformly positive elliptic operator and $Q_{0}$ corresponds to the $L^{2}$ duality, the operator $\mathbb{D}$ is basically a differential operator on $X_{+}$of the same order as $Q_{1}$. The symmetric operator $\mathbb{L}$ admits a resolution of identity consisting of bounded projections $\Pi_{\lambda}: X_{+} \rightarrow X_{+}$, $\lambda \in[0,1]$, where
(1) $\lim _{\lambda \rightarrow \lambda_{0}+} \Pi_{\lambda} u=\Pi_{\lambda_{0}} u$, for all $\lambda_{0} \in[0,1)$ and $u \in X_{+}$;
(2) $\Pi_{\lambda_{1}} \Pi_{\lambda_{2}}=\Pi_{\min \left\{\lambda_{1}, \lambda_{2}\right\}}$, for all $\lambda_{1,2} \in[0,1]$;
(3) $\left\langle\left(L_{+}+a Q_{0}^{+}\right) \Pi_{\lambda} u_{1}, u_{2}\right\rangle=\left\langle\left(L_{+}+a Q_{0}^{+}\right) u_{1}, \Pi_{\lambda} u_{2}\right\rangle$ for any $u \in X_{+}$and $\lambda \in[0,1] ;$
(4) $u=\int_{0}^{1} d \Pi_{\lambda} u, \mathbb{L} u=\int_{0}^{1} \lambda d \Pi_{\lambda} u$, for any $u \in X_{+}$.

Here, $\Pi_{1}=I$ and $\Pi_{0}=0$ since $L_{+}$is bounded and $0<L_{+}<L_{+}+a Q_{0}^{+}$as a quadratic form. Using this resolution of identity, we have the representations of $L_{+}$and $\|\cdot\|_{L_{+}, a}$

$$
\left\langle L_{+} u, v\right\rangle=\int_{0}^{1} \lambda d\left(\Pi_{\lambda} u, v\right)_{L_{+}, a},\|u\|_{L_{+}, a}^{2}=\int_{0}^{1} d\left\|\Pi_{\lambda} u\right\|_{L_{+}, a}^{2}, u, v \in X_{+}
$$

Let $\left(Y_{+},\|\cdot\|_{L_{+}}\right)$be the Hilbert space of the completion of $X_{+}$with respect to the inner product

$$
(u, v)_{L_{+}}=(\mathbb{L} u, v)_{L_{+}, a}=\left\langle L_{+} u, v\right\rangle=\langle L u, v\rangle=\int_{0}^{1} \lambda d\left(\Pi_{\lambda} u, v\right)_{L_{+}, a}, u, v \in X_{+}
$$

Therefore, $X_{+}$is densely embedded into $Y_{+}$through the embedding $i_{X_{+}}$. Using the above spectral integral representation of $\mathbb{L}$, one may extend $\Pi_{\lambda}$ to be bounded linear projections on $Y$ orthogonal with respect to $(\cdot, \cdot)_{L_{+}}$as well, satisfying $\left|\Pi_{\lambda}\right|_{Y} \leq 1$. Moreover, for $\lambda \in(0,1],\left(I-\Pi_{\lambda}\right) Y_{+} \subset X_{+}$and

$$
\begin{align*}
& \forall u \in X_{+},\left\|\Pi_{\lambda} u\right\|_{L_{+}} \leq \lambda\left\|\Pi_{\lambda} u\right\|_{L_{+}, a} \\
& \forall u \in Y_{+}, \lambda\left\|\left(I-\Pi_{\lambda}\right) u\right\|_{L_{+}, a} \leq\left\|\left(I-\Pi_{\lambda}\right) u\right\|_{L_{+}} \leq\left\|\left(I-\Pi_{\lambda}\right) u\right\|_{L_{+}, a} \tag{10.5}
\end{align*}
$$

where $I-\Pi_{\lambda}=\int_{(\lambda, 1]} d \Pi_{\lambda}$ is used.
As $Y_{+}$is defined as the completion of $X_{+}$with respect to the metric $(\mathbb{L} u, u)_{L_{+}, a}$, elements in $Y_{+}$are defined via Cauchy sequences in $X_{+}$with respect to this metric. This is rather inconvenient technically. Instead, we give an integral representation of elements in $Y_{+}$and some linear quantities on $Y_{+}$using $\Pi_{\lambda}$ and the following lemma.

Lemma 10.2. $\lim _{\lambda \rightarrow 0+}\left\|\Pi_{\lambda} u\right\|_{L_{+}}=0$ for any $u \in Y_{+}$.
Proof. For any $\epsilon>0$, there exists $v \in X_{+}$such that $\|u-v\|_{L_{+}}<\frac{\epsilon}{2}$. Since $\lim _{\lambda \rightarrow 0+} \Pi_{\lambda} v=\Pi_{0} v=0$ in $X_{+}$, there exists $\lambda_{0}>0$ such that $\left\|\Pi_{\lambda} v\right\|_{L_{+}, a}<\frac{\epsilon}{2}$ for any $\lambda \in\left(0, \lambda_{0}\right)$. Therefore, for any $\lambda \in\left(0, \lambda_{0}\right)$,

$$
\left\|\Pi_{\lambda} u\right\|_{L_{+}} \leq\left\|\Pi_{\lambda}(u-v)\right\|_{L_{+}}+\left\|\Pi_{\lambda} v\right\|_{L_{+}} \leq\|u-v\|_{L_{+}}+\lambda\left\|\Pi_{\lambda} v\right\|_{L_{+}, a} \leq \epsilon
$$

The lemma is proved.

Corollary 10.1. For any $u, v \in Y_{+}$, we have

$$
\begin{aligned}
& u=\int_{0}^{1} d \Pi_{\lambda} u=-\int_{0}^{1} d\left(I-\Pi_{\lambda}\right) u=-\lim _{\lambda \rightarrow 0+} \int_{\lambda}^{1} d\left(I-\Pi_{\lambda}\right) u \\
& \mathbb{L} u=-\int_{0}^{1} \lambda d\left(I-\Pi_{\lambda}\right) u,\|u\|_{L_{+}}^{2}=-\int_{0}^{1} \lambda d\left\|\left(I-\Pi_{\lambda}\right) u\right\|_{L_{+}, a}^{2} \\
& \left\langle L_{+} u, v\right\rangle=-\int_{0}^{1} \lambda d\left(\left(I-\Pi_{\lambda}\right) u, v\right)_{L_{+}, a}=\lim _{\lambda \rightarrow 0+}\left\langle L_{+}\left(I-\Pi_{\lambda}\right) u,\left(I-\Pi_{\lambda}\right) v\right\rangle \\
& Y_{+}^{*}=\left\{f=\left(L_{+}+a Q_{0}^{+}\right) u \mid u \in X_{+}\right. \\
& \left.\|f\|_{Y_{+}^{*}}^{2}=-\int_{0}^{1} \lambda^{-1} d\left\|\left(I-\Pi_{\lambda}\right) u\right\|_{L_{+}, a}^{2}<\infty\right\} \subset X_{+}^{*}
\end{aligned}
$$

Here, the first integral converges in the $\|\cdot\|_{L_{+}}$norm and the minus signs are due to the non-increasing monotonicity of $\left\|\left(I-\Pi_{\lambda}\right) u\right\|_{L_{+}, a}^{2}$. With $\left(I-\Pi_{\lambda}\right) u \in X_{+}$ for $\lambda \in(0,1]$, these integral representations are more convenient than the Cauchy sequence representations of elements in $Y_{+}$. In particular, for $f=\left(L_{+}+a Q_{0}^{+}\right) u \in$ $Y_{+}^{*}$ and $v \in Y_{+}$,

$$
\begin{aligned}
\langle f, v\rangle & =-\int_{0}^{1} d\left(\left(I-\Pi_{\lambda}\right) u, v\right)_{L_{+}, a}=\lim _{\lambda \rightarrow 0+}\left\langle\left(L_{+}+a Q_{0}^{+}\right)\left(I-\Pi_{\lambda}\right) u,\left(I-\Pi_{\lambda}\right) v\right\rangle \\
& \leq\|f\|_{Y_{+}^{*}}\|v\|_{L_{+}}
\end{aligned}
$$

Let

$$
\begin{equation*}
Y=X_{\leq 0} \oplus Y_{+}, \quad(u, v)_{Y}=\left(P_{\leq 0} u, P_{\leq 0} v\right)+\left(\left(I-P_{\leq 0}\right) u,\left(I-P_{\leq 0}\right) v\right)_{L_{+}} \tag{10.6}
\end{equation*}
$$

where, with slight abuse of notations, $P_{\leq 0}: Y \rightarrow X_{\leq 0}$ represents the projection operator with kernel $Y_{+}$. Clearly, $X$ is densely embedded into $Y$ and let $i_{X}$ denote the embedding.

The dual space $Y^{*}$ is densely embedded into $X^{*}$ through $i_{X}^{*}$ and thus can be viewed as a dense subspace of $X^{*}$. It is straightforward to see that $i_{X}^{*} Y^{*}=$ $\tilde{X}_{\leq 0}^{*} \oplus \tilde{Y}_{+}^{*}$,

$$
\langle f, v\rangle=\langle g, u\rangle=0, \forall u \in X_{\leq 0}, v \in Y_{+}, \quad f \in \tilde{X}_{\leq 0}^{*}, g \in \tilde{Y}_{+}^{*}
$$

and

$$
\begin{align*}
& \tilde{Y}_{+}^{*}=\tilde{X}_{+}^{*} \cap i_{X}^{*}\left(Y^{*}\right)=P_{+}^{*}\left\{f=\left(L_{+}+a Q_{0}^{+}\right) u \mid u \in X_{+}\right. \\
& \left.\|f\|_{Y_{+}^{*}}^{2}=-\int_{0}^{1} \lambda^{-1} d\left\|\left(I-\Pi_{\lambda}\right) u\right\|_{L_{+}, a}^{2}<\infty\right\} \subset X^{*} . \tag{10.7}
\end{align*}
$$

Operator $L$ is naturally extended as a bounded symmetric linear operator $L_{Y}$ : $Y \rightarrow Y^{*}$ by

$$
\begin{equation*}
\left\langle L_{Y} u, v\right\rangle=\left\langle L_{\leq 0} P_{\leq 0} u, P_{\leq 0} v\right\rangle+\left\langle L_{+}\left(I-P_{\leq 0}\right) u,\left(I-P_{\leq 0}\right) v\right\rangle \tag{10.8}
\end{equation*}
$$

where $L_{+}$on $Y_{+}$is computed by the formula given in Corollary 10.1.
From assumption ( $\mathbf{B 4}$ ) on the $L$-orthogonality of the decomposition

$$
X=X_{-} \oplus \operatorname{ker} L \oplus X_{+}=X_{\leq 0} \oplus X_{+}
$$

and Corollary 10.1, the operator $L_{Y}$ defined in the above satisfies $(\mathbf{H 2})$ on $Y$ with $\delta=1$ in (H2.b).

Operator $J_{Y}$. We define $J_{Y}: Y^{*} \supset D\left(J_{Y}\right) \rightarrow Y$ essentially as the restriction of $J$ on $Y^{*}$, namely,

$$
\begin{equation*}
J_{Y} \triangleq i_{X} \mathbb{J} Q_{0}^{-1} i_{X}^{*}: Y^{*} \supset D\left(J_{Y}\right) \rightarrow Y, \quad D\left(J_{Y}\right)=\left(i_{X}^{*}\right)^{-1} Q_{0}(X) \subset Y^{*} \tag{10.9}
\end{equation*}
$$

where we recall that $i_{X}: X \rightarrow Y$ is the embedding. Assumption (H3) is satisfied due to (B5) and (10.1). Therefore, to complete the proof of Theorem 2.7, it suffice to prove $J_{Y}^{*}=-J_{Y}$.

Lemma 10.3. It holds that $i_{X}^{*} D\left(J_{Y}\right)$ is dense in $X^{*}$ and

$$
i_{X}^{*} D\left(J_{Y}\right)=Q_{0}(X) \cap i_{X}^{*} Y^{*}=\tilde{X}_{\leq 0}^{*} \oplus P_{+}^{*}\left(L_{+}+a Q_{0}^{+}\right) X_{1+}
$$

where

$$
X_{1+}=\left\{u \in X_{+} \left\lvert\, \int_{0}^{1} \frac{-1}{\lambda(1-\lambda)^{2}} d\left\|\left(I-\Pi_{\lambda}\right) u\right\|^{2}<\infty\right.\right\} \subset X_{+}
$$

Proof. From $i_{X}^{*} Y^{*}=\tilde{X}_{\leq 0}^{*} \oplus \tilde{Y}_{+}^{*}$ and (B5), we can decompose

$$
i_{X}^{*} D\left(J_{Y}\right)=Q_{0}(X) \cap i_{X}^{*} Y^{*}=\tilde{X}_{\leq 0}^{*} \oplus\left(\tilde{Y}_{+}^{*} \cap Q_{0}(X)\right)
$$

As $\tilde{Y}_{+}^{*} \subset \tilde{X}_{+}^{*}$, we obtain from Lemma 10.1

$$
\begin{equation*}
i_{X}^{*} D\left(J_{Y}\right)=\tilde{X}_{\leq 0}^{*} \oplus\left(\tilde{Y}_{+}^{*} \cap P_{+}^{*} Q_{0}^{+}\left(X_{+}\right)\right) \tag{10.10}
\end{equation*}
$$

Recall (10.3) and we have $\left(L_{+}+a Q_{0}^{+}\right) u \in Q_{0}^{+}\left(X_{+}\right), u \in X_{+}$, if and only if $u \in D(\mathbb{D})$, which is equivalent to $u \in(I-\mathbb{L})\left(X_{+}\right)$according to (10.4). Therefore, we obtain that

$$
\left(L_{+}+a Q_{0}^{+}\right) u \in Q_{0}^{+}\left(X_{+}\right), u \in X_{+},
$$

if and only if

$$
-\int_{0}^{1}(1-\lambda)^{-2} d\left\|\left(I-\Pi_{\lambda}\right) u\right\|_{L_{+}, a}^{2}<\infty
$$

which can be seen from

$$
\begin{equation*}
\mathbb{D} u=\left(Q_{0}^{+}\right)^{-1}\left(L_{+}+a Q_{0}^{+}\right) u=-\frac{1}{a} \int_{0}^{1}(1-\lambda)^{-1} d\left(I-\Pi_{\lambda}\right) u \tag{10.11}
\end{equation*}
$$

The lemma follows immediately from this property and the characterization (10.7) of $\tilde{Y}_{+}^{*}$.

To prove $J_{Y}^{*}=-J_{Y}$, suppose $f \in D\left(J_{Y}^{*}\right)$ and $u=J_{Y}^{*} f$, namely, $f \in Y^{*}$ and $u \in Y$ satisfies

$$
\begin{equation*}
\left\langle f, J_{Y} g\right\rangle=\langle g, u\rangle, \quad \forall g \in D\left(J_{Y}\right) \tag{10.12}
\end{equation*}
$$

Firstly, for any $\epsilon \in\left(0, \frac{1}{2}\right)$, take

$$
g=-\left(i_{X}^{*}\right)^{-1} P_{+}^{*}\left(L_{+}+a Q_{0}^{+}\right) \int_{\epsilon}^{\frac{1}{2}} d\left(I-\Pi_{\lambda}\right) P_{+} u \in D\left(J_{Y}\right) \cap\left(i_{X}^{*}\right)^{-1} \tilde{X}_{+}^{*}
$$

where Lemma 10.3 is used. Equalities (10.12) and (10.11) imply

$$
\begin{aligned}
& -\int_{\epsilon}^{\frac{1}{2}} d\left\|\left(I-\Pi_{\lambda}\right) P_{+} u\right\|_{L_{+}, a}^{2}=\langle g, u\rangle=\left\langle f, J_{Y} g\right\rangle=\left\langle\mathbb{J}^{*} i_{X}^{*} f, Q_{0}^{-1} i_{X}^{*} g\right\rangle \\
= & -\left\langle\mathbb{J}^{*} i_{X}^{*} f, A \mathbb{D} \int_{\epsilon}^{\frac{1}{2}} d\left(I-\Pi_{\lambda}\right) P_{+} u\right\rangle=-\frac{1}{a}\left\langle\mathbb{J}^{*} i_{X}^{*} f, A \int_{\epsilon}^{\frac{1}{2}}(1-\lambda)^{-1} d\left(I-\Pi_{\lambda}\right) P_{+} u\right\rangle,
\end{aligned}
$$

where $A$ is defined in Lemma 10.1 and proved to be bounded. Since $\lambda \in\left(0, \frac{1}{2}\right]$, there exists $C>0$ such that

$$
-\int_{\epsilon}^{\frac{1}{2}} d\left\|\left(I-\Pi_{\lambda}\right) P_{+} u\right\|_{L_{+}, a}^{2} \leq C\|f\|_{Y^{*}}^{2}, \quad \forall \epsilon \in\left(0, \frac{1}{2}\right)
$$

Therefore, we obtain $u \in X$, or more precisely,

$$
u=i_{X} \tilde{u}, \quad \tilde{u} \in X
$$

For any $g \in D\left(J_{Y}\right)$, from (B1-2) we can compute

$$
\begin{aligned}
& \left\langle i_{X}^{*} f, \mathbb{J} Q_{0}^{-1} i_{X}^{*} g\right\rangle=\left\langle f, J_{Y} g\right\rangle=\langle g, u\rangle=\left\langle g, i_{X} \tilde{u}\right\rangle \\
= & \left\langle Q_{0} Q_{0}^{-1} i_{X}^{*} g, \tilde{u}\right\rangle=\left\langle Q_{0} \mathbb{J} Q_{0}^{-1} i_{X}^{*} g, \mathbb{J} \tilde{u}\right\rangle=\left\langle Q_{0} \mathbb{J} \tilde{u}, \mathbb{J} Q_{0}^{-1} i_{X}^{*} g\right\rangle .
\end{aligned}
$$

Since $\mathbb{J}$ is assumed to be isomorphic in (B2), $Q_{0}^{-1}$ is surjective, and $i_{X}^{*} D\left(J_{Y}\right)$ is dense in $X^{*}$ (Lemma 10.3), we obtain

$$
i_{X}^{*} f=Q_{0} \mathbb{J} \tilde{u} \in Q_{0}(X)
$$

Thus it follows from (B2) that

$$
J_{Y} f=i_{X} \mathbb{J}^{2} \tilde{u}=-u
$$

which implies $J_{Y}^{*} \subset-J_{Y}$. Again from (B2), it is easy to see that $J_{Y}$ is symmetric, namely, $J_{Y} \subset-J_{Y}^{*}$. Therefore, we complete the proof of Theorem 2.7.

## CHAPTER 11

## Hamiltonian PDE models

In this chapter, based on the above general theory, we study the stability issues of examples of Hamiltonian PDEs including several dispersive wave models, the 2D Euler equation for inviscid flows and a 2D nonlinear Schrödinger equations with nonzero conditions at infinity.

First, in Sections 11.1 to 11.3 , we study the stability/instability of traveling solitary and periodic wave solutions of several classes of equations modeling weakly nonlinear dispersive long waves. They include BBM, KDV, and good Boussinesq type equations. These equations respectively have the forms:

1. BBM type

$$
\begin{equation*}
\partial_{t} u+\partial_{x} u+\partial_{x} f(u)+\partial_{t} \mathcal{M} u=0 \tag{11.1}
\end{equation*}
$$

2. KDV type

$$
\begin{equation*}
\partial_{t} u+\partial_{x} f(u)-\partial_{x} \mathcal{M} u=0 \tag{11.2}
\end{equation*}
$$

3. good Boussinesq (gBou) type

$$
\begin{equation*}
\partial_{t}^{2} u-\partial_{x}^{2} u+\partial_{x}^{2} f(u)-\partial_{x}^{2} \mathcal{M} u=0 \tag{11.3}
\end{equation*}
$$

We follow the notations in [52]. Here, the pseudo-differential operator $\mathcal{M}$ is defined as

$$
\widehat{\mathcal{M} g}(\xi)=\alpha(\xi) \widehat{g}(\xi)
$$

where $\hat{g}$ is the Fourier transformation of $g$. We assume: i) $f$ is $C^{1}$ with $f(0)=$ $f^{\prime}(0)=0$, and $f(u) / u \rightarrow \infty$. ii) $a|\xi|^{m} \leq \alpha(\xi) \leq b|\xi|^{m}$ for large $\xi$, where $m>0$ and $a, b>0$. If $f(u)=u^{2}$ and $\mathcal{M}=-\partial_{x}^{2}$, the above equations recover the original BBM, KDV, and good Boussinesq equations, which have been used to model the propagation of water waves of long wavelengths and small amplitude.

### 11.1. Stability of Solitary waves of Long wave models

Consider the equations (11.1)-(11.3) with $(x, t) \in \mathbf{R} \times \mathbf{R}$. Up to a shift of a constant of the wave speed/symbol $\alpha(\xi)$, we can assume that $\sigma_{\text {ess }}(\mathcal{M}) \subset[0, \infty)$. Each of the equations (11.1)-(11.3) admits solitary-wave solutions of the form $u(x, t)=u_{c}(x-c t)$ for $c>1, c>0, c^{2}<1$ respectively, where $u_{c}(x) \rightarrow 0$ as $|x| \rightarrow \infty$. They satisfy the equations

$$
\begin{gather*}
\mathcal{M} u_{c}+\left(1-\frac{1}{c}\right) u_{c}-\frac{1}{c} f\left(u_{c}\right)=0,(\mathrm{BBM})  \tag{11.4}\\
\mathcal{M} u_{c}+c u_{c}-f\left(u_{c}\right)=0,(\mathrm{KDV})
\end{gather*}
$$

and

$$
\mathcal{M} u_{c}+\left(1-c^{2}\right) u_{c}-f\left(u_{c}\right)=0,(\text { gBou })
$$

respectively. We refer to the introduction of [52] and the book [2] for the literature on the existence of such solitary waves. Before stating the results, we introduce some notations. For BBM type equations (11.1), define the operator

$$
\begin{equation*}
\mathcal{L}_{0}=\mathcal{M}+\left(1-\frac{1}{c}\right)-\frac{1}{c} f^{\prime}\left(u_{c}\right): H^{m} \rightarrow L^{2} \tag{11.5}
\end{equation*}
$$

and the momentum

$$
\begin{equation*}
P(c)=\frac{1}{2} \int u_{c}(\mathcal{M}+1) u_{c} . \tag{11.6}
\end{equation*}
$$

For KDV type equations (11.2), define

$$
\begin{equation*}
\mathcal{L}_{0}:=\mathcal{M}+c-f^{\prime}\left(u_{c}\right), \quad P(c)=\frac{1}{2} \int u_{c}^{2} \tag{11.7}
\end{equation*}
$$

For good Boussinesq type equations (11.3), define

$$
\begin{equation*}
\mathcal{L}_{0}:=\mathcal{M}+1-c^{2}-f^{\prime}\left(u_{c}\right), P(c)=-c \int u_{c}^{2} \tag{11.8}
\end{equation*}
$$

Denote by $n^{-}\left(\mathcal{L}_{0}\right)$ the number (counting multiplicity) of negative eigenvalues of the operators $\mathcal{L}_{0}$.

The linearizations of (11.1)-(11.3) in the traveling frame $(x-c t, t)$ are

$$
\begin{gather*}
\left(\partial_{t}-c \partial_{x}\right)(u+\mathcal{M} u)+\partial_{x}\left(u+f^{\prime}\left(u_{c}\right) u\right)=0,(\mathrm{BBM})  \tag{11.9}\\
\left(\partial_{t}-c \partial_{x}\right) u+\partial_{x}\left(f^{\prime}\left(u_{c}\right) u-\mathcal{M} u\right)=0,(\mathrm{KDV}) \tag{11.10}
\end{gather*}
$$

and

$$
\begin{equation*}
\left(\partial_{t}-c \partial_{x}\right)^{2} u-\partial_{x}^{2}\left(u-f^{\prime}\left(u_{c}\right) u+\mathcal{M} u\right)=0,(\text { gBou }) \tag{11.11}
\end{equation*}
$$

respectively. We consider the Hamiltonian structures of these equations.
For BBM type equations, (11.9) can be written as $\partial_{t} u=J L u$, where $J=$ $c \partial_{x}(1+\mathcal{M})^{-1}$ and $L=\mathcal{L}_{0}$ is defined in (11.5). By differentiating (11.4) in $x$ and $c$, we have $\mathcal{L}_{0} u_{c, x}=0$ and

$$
\mathcal{L}_{0} \partial_{c} u_{c}=-\frac{1}{c}(1+\mathcal{M}) u_{c}
$$

which implies that $J \mathcal{L}_{0} \partial_{c} u_{c}=-u_{c, x}$ and $\left\langle\mathcal{L}_{0} \partial_{c} u_{c}, \partial_{c} u_{c}\right\rangle=-\frac{1}{c} d P / d c$.
For KDV type equations, (11.10) is written as $\partial_{t} u=J \mathcal{L}_{0} u$, where $J=\partial_{x}$ and $L=\mathcal{L}_{0}$ is defined in (11.7). Similarly, $\mathcal{L}_{0} u_{c, x}=0, \mathcal{L}_{0} \partial_{c} u_{c}=-u_{c}$, and

$$
\begin{equation*}
J \mathcal{L}_{0} \partial_{c} u_{c}=-u_{c, x},\left\langle\mathcal{L}_{0} \partial_{c} u_{c}, \partial_{c} u_{c}\right\rangle=-d P / d c \tag{11.12}
\end{equation*}
$$

For good Boussinesq type equations, we write (11.11) as a first order system. Let $\left(\partial_{t}-c \partial_{x}\right) u=v_{x}$, then

$$
\left(\partial_{t}-c \partial_{x}\right) v=\partial_{x}\left(\mathcal{M}+1-f^{\prime}\left(u_{c}\right)\right) u=\partial_{x}\left(\mathcal{L}_{0}+c^{2}\right) u
$$

Thus

$$
\partial_{t}\binom{u}{v}=J L\binom{u}{v}
$$

with

$$
J=\left(\begin{array}{cc}
0 & \partial_{x}  \tag{11.13}\\
\partial_{x} & 0
\end{array}\right), \quad L=\left(\begin{array}{cc}
\mathcal{L}_{0}+c^{2} & c \\
c & 1
\end{array}\right)
$$

We have

$$
\operatorname{ker} L=\left\{(u,-c u) \mid u \in \operatorname{ker} \mathcal{L}_{0}\right\}
$$

Since

$$
\begin{equation*}
\left\langle L\binom{u}{v},\binom{u}{v}\right\rangle=\left\langle\mathcal{L}_{0} u, u\right\rangle+\int(v+c u)^{2} \tag{11.14}
\end{equation*}
$$

so $n^{-}(L)=n^{-}\left(\mathcal{L}_{0}\right)$. Similarly as in BBM and KDV types, we have

$$
\begin{gathered}
L\binom{u_{c, x}}{-c u_{c, x}}=0, L\binom{-\partial_{c} u_{c}}{c \partial_{c} u_{c}+u_{c}}=\binom{-c u_{c}}{u_{c}}, \\
J L\binom{-\partial_{c} u_{c}}{c \partial_{c} u_{c}+u_{c}}=\binom{u_{c, x}}{-c u_{c, x}}
\end{gathered}
$$

and

$$
\left\langle L\binom{-\partial_{c} u_{c}}{c \partial_{c} u_{c}+u_{c}},\binom{-\partial_{c} u_{c}}{c \partial_{c} u_{c}+u_{c}}\right\rangle=-d P / d c
$$

where $P$ is defined in (11.8). For all three cases, we have $\sigma_{\text {ess }}\left(\mathcal{L}_{0}\right) \subset\left[\delta_{0}, \infty\right)$ for some $\delta_{0}>0$. So the quadratic form $\left\langle\mathcal{L}_{0} \cdot, \cdot\right\rangle$ is positive definite on $H^{\frac{m}{2}}$ in a finite codimensional space. This along with Remark 2.3 shows that the quadratic form $\left\langle\mathcal{L}_{0} \cdot, \cdot\right\rangle$ in the Hamiltonian formulation of BBM and KDV type equations satisfies the assumption (H1-3) in the general framework with $X=H^{\frac{m}{2}}$. By (11.14), the quadratic form $\langle L \cdot, \cdot\rangle$ in the Hamiltonian formulation (11.13) of good Boussinesq type equations also satisfies (H1-3) in the space $\left(u, \partial_{t} u\right) \in X=H^{\frac{m}{2}} \times L^{2}$. Thus by Theorems 2.2, 2.3, and Corollary 2.2, we get the following results.

Theorem 11.1. Consider the linearized equations (11.9)-(11.11) at solitary waves $u_{c}(x-c t)$ of equations (11.1)-(11.3). Then: (i) The following index formula holds

$$
\begin{equation*}
k_{r}+2 k_{c}+2 k_{i}^{\leq 0}+k_{0}^{\leq 0}=n^{-}\left(\mathcal{L}_{0}\right) . \tag{11.15}
\end{equation*}
$$

(ii) The linear exponential trichotomy holds in the space $H^{\frac{m}{2}}$ for the linearized equations (11.9) and (11.10), and in $H^{\frac{m}{2}} \times L^{2}$ for (11.11).
(iii) When $d P / d c \geq 0$, we have $k_{0}^{\leq 0} \geq 1$. Moreover, if $\operatorname{ker} \mathcal{L}_{0}=\operatorname{span}\left\{u_{c, x}\right\}$, then

$$
k_{0}^{\leq 0}=\left\{\begin{array}{ll}
1 & \text { if } d P / d c>0 \\
0 & \text { if } d P / d c<0
\end{array} .\right.
$$

Corollary 11.1. (i) When $d P / d c \geq 0$ and $n^{-}\left(\mathcal{L}_{0}\right) \leq 1$, the spectral stability holds true.
(ii) If $\operatorname{ker} \mathcal{L}_{0}=\operatorname{span}\left\{u_{c, x}\right\}$, then there is linear instability when $n^{-}\left(\mathcal{L}_{0}\right)$ is even and $d P / d c>0$ or $n^{-}\left(\mathcal{L}_{0}\right)$ is odd and $d P / d c<0$.

In particular, when $\mathcal{M}=-\partial_{x}^{2}$, by the fact that $u_{c, x}$ changes sign exactly once and the Sturm-Liouville theory, we have $\operatorname{ker} \mathcal{L}_{0}=\operatorname{span}\left\{u_{c, x}\right\}$ and $n^{-}\left(\mathcal{L}_{0}\right)=1$. Thus, we have

Corollary 11.2. When $\mathcal{M}=-\partial_{x}^{2}$ and $d P / d c<0$, for the linearized equations (11.1)-(11.3), we have $k_{r}=1$ and $k_{c}=k_{i}^{-}=0$. In particular, on the center space $E^{c}$ as given in Theorem 2.2, we have

$$
\begin{equation*}
\left.\langle L \cdot, \cdot\rangle\right|_{E^{c} \cap\left\{u_{c, x}\right\}^{\perp}} \geq \delta_{0}>0 \tag{11.16}
\end{equation*}
$$

The stability and instability of solitary waves of dispersive models had been studied a lot in the literature. Assume $\operatorname{ker} \mathcal{L}_{0}=\left\{u_{c, x}\right\}, n^{-}\left(\mathcal{L}_{0}\right)=1$, then when $d P / d c>0$, the orbital stability of traveling solitary waves was proved (e.g. [5] [29] [10]) by using the method of Lyapunov functionals. When $d P / d c<0$, the nonlinear
instability was proved in [9] [70] for generalized BBM and KDV equations, and in [58] for good Boussinesq equation. The instability proof in these papers was by contradiction argument which bypassed the linearized equation. The existence of unstable eigenvalues when $d P / d c>0$ was proved in [63] for KDV and BBM equations. In [52], an instability criterion as in Corollary 11.1 (ii) was proved for KDV and BBM type equations. In [63] and [52], an instability criterion was also given for the regularized Boussinesq equation which takes an indefinite Hamiltonian form (i.e. $\left.n^{-}(L)=\infty\right)$ and is therefore not included in the framework of this paper. Recently, in [46] and [65], an instability index theorem similar to (11.15) was given for KDV and BBM type equations under the assumption that $\operatorname{dim} \operatorname{ker} \mathcal{L}_{0}=1$ and $d P / d c \neq 0$. The proof of [46] [65] was by using ad-hoc arguments to transform the eigenvalue problem $\partial_{x} \mathcal{L}_{0} u=\lambda u$ to another Hamiltonian form with a symplectic operator which has a bounded inverse. The linear instability of solitary waves of good Boussinesq equation $\left(\mathcal{M}=-\partial_{x}^{2}\right)$ was studied in [1] by Evans function and in [68] by using quadratic operator pencils. The index formula (11.15) for the good Boussinesq type equations appears to be new.

Besides giving a more unified and general index formula for linear instability, Theorem 11.1 also gives the exponential trichotomy for $e^{t J \mathcal{L}_{0}}$, which is an important step for constructing invariant (stable, unstable and center) manifolds near the translation orbits of $u_{c}$. Moreover, when $\operatorname{ker} \mathcal{L}_{0}=\operatorname{span}\left\{u_{c, x}\right\}$ and $n^{-}\left(\mathcal{L}_{0}\right)=$ 1 , there exists a pair of stable and unstable eigenvalues and $\mathcal{L}_{0}$ is positive on the codimension two center space modulo the translation kernel. This positivity property has an important implication for the center manifolds once constructed. For example, in [37], the invariant (stable, unstable and center) manifolds were constructed near the orbits of unstable solitary waves of generalized KDV equation in the energy space. More precisely, there exist 1-d stable and unstable manifolds and co-dimension two center manifold near the translation orbits of unstable solitary waves. These invariant manifolds give a complete description of local dynamics near unstable traveling wave orbits. The positivity estimate (11.16) on the center subspace implies that on the codimension two center manifold, the solitary wave $u_{c}$ is orbitally stable, which in turn also leads to the local uniqueness of the center manifold. Any initial data not lying on the center manifold will leave the orbit neighborhood of unstable traveling waves exponentially fast.

### 11.2. Stability of periodic traveling waves

Consider the equations (11.1)-(11.3) in the periodic case. For convenience, we assume the period is $2 \pi$, that is, $(x, t) \in \mathbf{S}^{1} \times \mathbf{R}$. A periodic traveling wave is of the form $u(x, t)=u_{c, a}(x-c t)$, where $u_{c, a}$ satisfies the equations

$$
\begin{gather*}
\mathcal{M} u_{c, a}+\left(1-\frac{1}{c}\right) u_{c, a}-\frac{1}{c} f\left(u_{c, a}\right)=a,(\mathrm{BBM})  \tag{11.17}\\
\mathcal{M} u_{c, a}+c u_{c, a}-f\left(u_{c, a}\right)=a,(\mathrm{KDV}) \tag{11.18}
\end{gather*}
$$

and

$$
\begin{equation*}
\mathcal{M} u_{c, a}+\left(1-c^{2}\right) u_{c, a}-f\left(u_{c, a}\right)=a,(\mathrm{gBou}) \tag{11.19}
\end{equation*}
$$

for some constant $a$. In this section, we consider the perturbations of the same period $2 \pi$ (i.e. co-periodic perturbations) and leave the case of different periods to the next section. The linearized equations in the traveling frame $(x-c t, t)$ near
traveling waves $u_{c, a}$ take the same form (11.9)-(11.11). Their Hamiltonian structures are formally the same as in the case of solitary waves. However, the operator $J$ has rather different spectral properties in the periodic case. More precisely, for solitary waves the symplectic operators $J$, which is $c \partial_{x}(1+\mathcal{M})^{-1}$ for BBM, $\partial_{x}$ for KDV and $\left(\begin{array}{cc}0 & \partial_{x} \\ \partial_{x} & 0\end{array}\right)$ for good Boussinesq, has no kernel in $L^{2}(\mathbf{R})$. But for the periodic case, $J$ has nontrivial kernel in $X^{*}$. Indeed, $\operatorname{ker} J=\operatorname{span}\{1\}$ for BBM and KDV, and

$$
\text { ker } J=\operatorname{span}\left\{\vec{e}_{1}, \vec{e}_{2}\right\}=\operatorname{span}\left\{\binom{1}{0},\binom{0}{1}\right\}
$$

for good Boussinesq. This degeneracy of $J$ leads to the extra free parameter $a$ in traveling waves.

We now discuss the consequential changes in the index formula induced by the nontrivial kernel of $J$. For BBM type equations, define the operator $\mathcal{L}_{0}$ : $H^{m}\left(\mathbf{S}^{1}\right) \rightarrow L^{2}\left(\mathbf{S}^{1}\right)$ and the momentum $P$ as in (11.5) and (11.6). Differentiating (11.17), we obtain

$$
R\left(\mathcal{L}_{0}\right) \ni \mathcal{L}_{0} \partial_{a} u_{c, a}=1
$$

Let

$$
\begin{equation*}
U_{c, a}=\partial_{a} u_{c, a}, d_{1}=\int_{\mathbf{S}^{1}} U_{c, a} d x, N=\int_{\mathbf{S}^{1}} u_{c, a} d x \text { (total mass). } \tag{11.20}
\end{equation*}
$$

We have $\mathcal{L}_{0} \partial_{x} u_{c, a}=0$ and from differentiating (11.17)

$$
\mathcal{L}_{0} \partial_{c} u_{c, a}=-\frac{1}{c}(1+\mathcal{M}) u_{c, a}+\frac{a}{c}
$$

and thus $J \mathcal{L}_{0} \partial_{c} u_{c, a}=-\partial_{x} u_{c, a}$. Denote

$$
D=\left(\begin{array}{cc}
\left\langle\mathcal{L}_{0} U_{c, a}, U_{c, a}\right\rangle & \left\langle\mathcal{L}_{0} U_{c, a}, \partial_{c} u_{c, a}\right\rangle  \tag{11.21}\\
\left\langle\mathcal{L}_{0} U_{c, a}, \partial_{c} u_{c, a}\right\rangle & \left\langle\mathcal{L}_{0} \partial_{c} u_{c, a}, \partial_{c} u_{c, a}\right\rangle
\end{array}\right)=\left(\begin{array}{cc}
d_{1} & N^{\prime}(c) \\
N^{\prime}(c) & -\frac{1}{c} d P / d c+\frac{a}{c} N^{\prime}(c)
\end{array}\right)
$$

that is, the matrix for $\langle L \cdot, \cdot\rangle$ on $\operatorname{span}\left\{U_{c, a}, \partial_{c} u_{c, a}\right\} \subset g \operatorname{ker}\left(J \mathcal{L}_{0}\right)$. Denote $n^{\leq 0}(D)$ to be the number of non-positive eigenvalues of $D$.

For KDV type equations, similarly, $\mathcal{L}_{0} \partial_{x} u_{c, a}=0$,

$$
\mathcal{L}_{0} \partial_{c} u_{c, a}=-u_{c, a}, J \mathcal{L}_{0} \partial_{c} u_{c, a}=-\partial_{x} u_{c, a}, \mathcal{L}_{0} \partial_{a} u_{c, a}=1
$$

and we define $U_{c, a}, d_{1}, N, D, n^{\leq 0}(D)$ etc. as in (11.20) and (11.21).
For good Boussinesq type equations, still define $U_{c, a}, d_{1}, N$ as in (11.20). Let

$$
\begin{equation*}
\vec{U}_{1}=\binom{U_{c, a}}{-c U_{c, a}}, \vec{U}_{2}=\binom{-c U_{c, a}}{1+c^{2} U_{c, a}}, \vec{U}_{3}=\binom{-\partial_{c} u_{c, a}}{c \partial_{c} u_{c, a}+u_{c, a}} \tag{11.22}
\end{equation*}
$$

then

$$
L \vec{U}_{1}=\vec{e}_{1}, \quad L \vec{U}_{2}=\vec{e}_{2}, \quad L \vec{U}_{3}=\binom{-c u_{c, a}}{u_{c, a}}
$$

Define the matrix $D$ of $\langle L \cdot, \cdot\rangle$ on the space spanned by $\left\{\vec{U}_{1}, \vec{U}_{2}, \vec{U}_{3}\right\}$, that is,

$$
\begin{align*}
D & =\left(\begin{array}{ccc}
\left\langle L \vec{U}_{1}, \vec{U}_{1}\right\rangle & \left.\left\langle\begin{array}{l}
\left.L \vec{U}_{1}, \vec{U}_{2}\right\rangle
\end{array}\right\rangle \begin{array}{l}
\left\langle\vec{U}_{1}, \vec{U}_{3}\right\rangle \\
\left\langle\vec{U}_{1}, \vec{U}_{2}\right. \\
\left\langle\vec{U}_{1}, \vec{U}_{3}\right.
\end{array}\right\rangle & \left\langle\begin{array}{l}
L \vec{U}_{2}, \vec{U}_{2}
\end{array}\right\rangle\left\langle\begin{array}{l}
\left\langle\vec{U}_{3}, \vec{U}_{2}\right. \\
\left\langle\vec{U}_{2}\right\rangle
\end{array}\right\rangle \\
\left\langle\vec{U}_{3}, \vec{U}_{3}\right\rangle
\end{array}\right)  \tag{11.23}\\
& =\left(\begin{array}{ccc}
c_{1} & -c d_{1} \\
-d_{1} & \int_{\mathbf{S}^{1}}\left(1+c^{2} U_{c, a}\right) d x & c N^{\prime}(c)+N(c) \\
-N^{\prime}(c) & c N^{\prime}(c)+N(c) & -P^{\prime}(c)
\end{array}\right) .
\end{align*}
$$

Again $n^{\leq 0}(D)$ denotes the number of non-positive eigenvalues of $D$.
Since in the periodic case, the operator $\mathcal{L}_{0}$ has only discrete spectrum which tends to $+\infty$, it is easy to verify that assumptions $(\mathbf{H} 1-3)$ are satisfied in $X=H^{\frac{m}{2}}$ for BBM and KdV type equations and $X=H^{\frac{m}{2}} \times L^{2}$ for good Boussinesq type equations. Thus similar to Theorem 11.1, we have

Theorem 11.2. Consider the linearized equations (11.9)-(11.11) near periodic waves $u_{c}(x-c t)$ of equations (11.1)-(11.3). Then: (i) the following index formula holds

$$
k_{r}+2 k_{c}+2 k_{i}^{-}+k_{0}^{-}=n^{-}\left(\mathcal{L}_{0}\right)
$$

(ii) the linear exponential trichotomy is true in the space $H^{\frac{m}{2}}\left(\mathbf{S}^{1}\right)$ for the linearized equations (11.9) and (11.10), and in $H^{\frac{m}{2}}\left(\mathbf{S}^{1}\right) \times L^{2}\left(\mathbf{S}^{1}\right)$ for (11.11). (iii) $k_{0}^{-} \geq$ $n^{\leq 0}(D)$, the number of non-positive eigenvalues of the matrix $D$ defined in (11.21), (11.22) and (11.23). Moreover, when $\operatorname{ker} \mathcal{L}_{0}=\left\{\partial_{x} u_{c, a}\right\}$ and $D$ is nonsingular, $k_{0}^{-}=n^{-}(D)$ (the number of negative eigenvalues of $D$ ) and we have

$$
\begin{equation*}
k_{r}+2 k_{c}+2 k_{i}^{-}=n^{-}\left(\mathcal{L}_{0}\right)-n^{-}(D) . \tag{11.24}
\end{equation*}
$$

As corollaries, we have from Proposition 2.7 and Remark 2.14 the following linear stability/instability conditions.

Corollary 11.3. (i) If $n^{\leq 0}(D) \geq n^{-}\left(\mathcal{L}_{0}\right)$, then the spectral stability holds.
(ii) If $\operatorname{ker} \mathcal{L}_{0}=\operatorname{span}\left\{\partial_{x} u_{c, a}\right\}, D$ is nonsingular and $n^{-}\left(\mathcal{L}_{0}\right)-n^{-}(D)$ is odd, then there is linear instability.

When $n^{-}\left(\mathcal{L}_{0}\right)=n^{-}(D)$, nonlinear orbital stability holds for (11.1)-(11.3) as well. More precisely, we have

Proposition 11.1. When ker $\mathcal{L}_{0}=\operatorname{span}\left\{\partial_{x} u_{c, a}\right\}, D$ is nonsingular and $n^{-}\left(\mathcal{L}_{0}\right)=$ $n^{-}(D)$, then there is orbital stability in $X$ of the traveling waves $u_{c, a}(x-c t)$ of equations (11.1)-(11.3) for perturbations of the same period.

Proof. Here we sketch the proof based on the standard Lyapunov functional method (e.g. [29], [30]). Consider the KDV type equation (11.2). It has three invariants: (1) energy $E(u)=\int\left[\frac{1}{2} u \mathcal{M} u-F(u)\right] d x$, with $F(u)=\int_{0}^{u} f\left(u^{\prime}\right) d u^{\prime}$; (2) momentum $P(u)=\frac{1}{2} \int u^{2} d x$ and (3) total mass $N(u)=\int u d x$. Define the invariant

$$
I(u)=E(u)+c P(u)-a N(u)
$$

then $I^{\prime}\left(u_{c, a}\right)=0$ if and only if $u_{c, a}$ is a traveling wave solution satisfying (11.18). So

$$
\begin{equation*}
I(u)-I\left(u_{c, a}\right)=\left\langle\mathcal{L}_{0} \delta u, \delta u\right\rangle+O\left(\|\delta u\|^{3}\right), \text { where } \delta u=u-u_{c, a} \tag{11.25}
\end{equation*}
$$

Denote

$$
X_{1}=\left\{\left.u \in H^{\frac{m}{2}} \right\rvert\,\left\langle u, \mathcal{L}_{0} U_{c, a}\right\rangle=\left\langle u, \mathcal{L}_{0} \partial_{c} u_{c, a}\right\rangle=0\right\}
$$

to be the orthogonal complement of $X_{2}=\operatorname{span}\left\{U_{c, a}, \partial_{c} u_{c, a}\right\}$ in $\left\langle\mathcal{L}_{0} \cdot, \cdot\right\rangle$. Since

$$
\mathcal{L}_{0} U_{c, a}=N^{\prime}\left(u_{c, a}\right), \quad \mathcal{L}_{0} \partial_{c} u_{c, a}=P^{\prime}\left(u_{c, a}\right)
$$

$X_{1}$ is the tangent space of the intersection of the level surfaces of the conserved momentum $P$ and mass $N$. With $D$ assumed to be nonsingular, $X_{2}$ roughly represents the gradient directions of $P$ and $N$ and thus $X=X_{1} \oplus X_{2}$. Moreover, we have $\left.\left\langle\mathcal{L}_{0} \cdot, \cdot\right\rangle\right|_{X_{1}} \geq 0$ since

$$
n^{-}\left(\left.\mathcal{L}_{0}\right|_{X_{1}}\right)=n^{-}\left(\mathcal{L}_{0}\right)-n^{-}\left(\left.\mathcal{L}_{0}\right|_{X_{2}}\right)=n^{-}\left(\mathcal{L}_{0}\right)-n^{-}(D)=0 .
$$

We further decompose

$$
X_{1}=Y \oplus \operatorname{span}\left\{\partial_{x} u_{c, a}\right\}, \quad \text { where } Y=\left\{u \in X_{1} \mid\left(u, \partial_{x} u_{c, a}\right)_{X}=0\right\}
$$

Since $\operatorname{ker} \mathcal{L}_{0}=\operatorname{span}\left\{\partial_{x} u_{c, a}\right\}$, there exists $c_{0}>0$ such that $\left\langle\mathcal{L}_{0} \delta u, \delta u\right\rangle \geq c_{0}\|\delta u\|^{2}$ for any $\delta u \in Y$.

Suppose $u(t)$ is solution with $u(0)$ close to $u_{c, a}$ and $h(t) \in \mathbf{S}^{1}$ satisfies

$$
\left\|u-u_{c, a}(\cdot-h)\right\|=\min _{y \in \mathbf{S}^{1}}\left\|u-u_{c, a}(\cdot-y)\right\|
$$

then $w(t)=u(t)-u_{c, a}(\cdot-h(t)) \in Y \oplus X_{2}$. By using the conservation of $P$ and $N$ to control the $X_{2}$ components of $w(t)$, and the uniform positivity of $\mathcal{L}_{0}$ on $Y$ and (11.25) to control the $Y$ component, we obtain the orbital stability. More details of such arguments can be found for example in [29] [53].

For BBM type equations, the Lyapunov functional is

$$
I(u)=c P(u)-E(u)-c a N(u),
$$

where the energy functional $E(u)=\int\left(\frac{1}{2} u^{2}+F(u)\right) d x$ and $P(u)$ is defined in (11.6). The rest of the proof is the same as in the KDV case. For good Boussinesq type equations (11.3), we write it as a first order Hamiltonian system

$$
\partial_{t}\binom{u}{v}=J \nabla E(u, v)
$$

where $J=\left(\begin{array}{cc}0 & \partial_{x} \\ \partial_{x} & 0\end{array}\right)$ and the energy functional

$$
E(u, v)=\frac{1}{2}(\mathcal{M} u, u)+\int\left(\frac{1}{2} v^{2}+\frac{1}{2} u^{2}-F(u)\right) d x
$$

For the traveling wave solution $\left(u_{c, a}(x-c t), v_{c, a}(x-c t)\right), u_{c, a}$ satisfies (11.17) and $v_{c, a}=-c u_{c, a}$. Let $\vec{u}=(u, v)^{T}$ and construct the Lyapunov functional

$$
I(\vec{u})=E(\vec{u})+c P(\vec{u})-a N_{1}(\vec{u}),
$$

where

$$
P(\vec{u})=\int u v d x, N_{1}(\vec{u})=\int u d x, N_{2}(\vec{u})=\int v d x
$$

Then $I^{\prime}\left(\vec{u}_{c, a}\right)=0$. The rest of the proof is the same.

Compared with solitary waves, the periodic traveling waves have richer structures. They consist of a three parameter (period $T$, speed $c$, and integration constant $a$ ) family of solutions and different type of perturbations (co-periodic, multiple periodic, localized etc.) can be considered. In recent years, there have been lots of works on stability/instability of periodic traveling waves of dispersive PDEs. For co-periodic perturbations (i.e. of the same period), the nonlinear orbital stability were proved for various dispersive models (e.g. [2] [3] [41] [34] [11] [8]) by using Liapunov functionals. These stability results were proved for the cases when $\operatorname{dim} \operatorname{ker}\left(\mathcal{L}_{0}\right)=1$ and $n^{-}\left(\mathcal{L}_{0}\right)=n^{-}(D)$ as in Proposition 11.1. An instability index formula similar to (11.24) was proved for KDV type equations ([32] [43] [13]). In these papers, some conditions (e.g. Assumption 2.1 in [32] and Assumption 3 in [43]) were imposed to ensure that the generalized eigenvectors of $J \mathcal{L}_{0}$ form a basis of $X$. These assumptions can be checked for the case $\mathcal{M}=-\partial_{x}^{2}$. In Theorem 11.2, we do not need such assumptions on the completion of generalized eigenspaces of $J \mathcal{L}_{0}$ and therefore we can get the index formula for very general nonlocal operators $\mathcal{M}$. In [11], an index formula was proved for periodic traveling waves of good Boussinesq equation $\left(\mathcal{M}=-\partial_{x}^{2}\right)$ by using the theory of quadratic operator pencils. In [8], a parity instability criterion (as in Corollary 11.3 (ii)) was proved for periodic waves of several Hamiltonian PDEs including generalized KDV equations by using Evans functions.

Besides providing a unified way to get instability index formula and the stability criterion, we could also use the exponential trichotomy of $e^{J \mathcal{L}_{0}}$ in Theorem 11.2 to construct invariant manifolds near the orbit of unstable periodic traveling waves. Moreover, as in the case of solitary waves, when $\operatorname{dim} \operatorname{ker}\left(\mathcal{L}_{0}\right)=1, D$ is nonsingular and $k_{i}^{-}=0$, we have orbital stability and local uniqueness of the center manifolds once constructed.

### 11.3. Modulational Instability of periodic traveling waves

Consider periodic traveling waves $u_{c, a}(x-c t)$ studied in Section 11.2. Assume the conditions in Proposition 11.1, so that $u_{c, a}$ is orbitally stable under perturbations of the same period. In this section, we consider modulational instability of periodic traveling waves, under perturbations of different period or even localized perturbations. The modulational instability, also called Benjamin-Feir or side-band instability in the literature, is a very important instability mechanism in lots of dispersive and fluid models. Again, we assume the minimal period of the traveling wave $u_{c, a}$ is $2 \pi$. We focus on KDV type equations (11.2), and the consideration for BBM and good-Boussinesq type equations is similar. We assume the Fourier symbol $\alpha(\xi)$ of the operator $\mathcal{M}$ is even, so that $\mathcal{M}$ is a real operator. Based on the standard Floquet-Bloch theory, we seek bounded eigenfunction $\phi(x)$ of the linearized operator $J \mathcal{L}_{0}$ in the form of $\phi(x)=e^{i k x} v_{k}(x)$, where $k \in \mathbf{R}$ is a parameter and $v_{k} \in L^{2}\left(\mathbf{S}^{1}\right)$. Recall that $J=\partial_{x}$ and $\mathcal{L}_{0}:=\mathcal{M}+c-f^{\prime}\left(u_{c, a}\right)$. It leads us to the one-parameter family of eigenvalue problems

$$
J \mathcal{L}_{0} e^{i k x} v_{k}(x)=\lambda(k) e^{i k x} v_{k}(x)
$$

or equivalently $\mathcal{J}_{k} \mathcal{L}_{k} v_{k}=\lambda(k) v_{k}$, where

$$
\begin{equation*}
\mathcal{J}_{k}=\partial_{x}+i k, \mathcal{L}_{k}=\mathcal{M}_{k}+c-f^{\prime}\left(u_{c, a}\right) . \tag{11.26}
\end{equation*}
$$

Here, $\mathcal{M}_{k}$ is the Fourier multiplier operator with the symbol $\alpha(\xi+k)$. We say that $u_{c, a}$ is linearly modulationally unstable if there exists $k \notin \mathbf{Z}$ such that the operator $\mathcal{J}_{k} \mathcal{L}_{k}$ has an unstable eigenvalue $\lambda(k)$ with $\operatorname{Re} \lambda(k)>0$ in the space $L^{2}\left(\mathbf{S}^{1}\right)$.

Since $\mathcal{J}_{k}$ and $\mathcal{L}_{k}$ are complex operators, we first reformulate the problem in terms of real operators to use the general theory in this paper. Consider

$$
\begin{equation*}
\phi(x)=\cos (k x) u_{1}(x)+\sin (k x) u_{2}(x), \tag{11.27}
\end{equation*}
$$

where $u_{1}, u_{2} \in L^{2}\left(\mathbf{S}^{1}\right)$ are real functions. By definition,

$$
\mathcal{M}\left(e^{i k x} u(x)\right)=e^{i k x} \mathcal{M}_{k} u
$$

We decompose

$$
\mathcal{M}_{k}=\mathcal{M}_{k}^{e}+i \mathcal{M}_{k}^{o}, \quad \mathcal{M}_{-k}=\mathcal{M}_{k}^{e}-i \mathcal{M}_{k}^{o}
$$

where $\mathcal{M}_{k}^{e}, \mathcal{M}_{k}^{o}$ are operators with Fourier multipliers

$$
\alpha_{k}^{e}(\xi)=\frac{1}{2}(\alpha(\xi+k)+\alpha(\xi-k))
$$

and

$$
\alpha_{k}^{o}(\xi)=-\frac{i}{2}(\alpha(\xi+k)-\alpha(\xi-k))
$$

Then $\mathcal{M}_{k}^{e}, \mathcal{M}_{k}^{o}$ are self-adjoint and skew-adjoint respectively. Since $\overline{\alpha_{k}^{e, o}(\xi)}=$ $\alpha_{k}^{e, o}(-\xi), \mathcal{M}_{k}^{e}$ and $\mathcal{M}_{k}^{o}$ map real functions to real. In particular, for $\mathcal{M}=-\partial_{x}^{2}$, we have $\mathcal{M}_{k}^{e}=-\partial_{x}^{2}+k^{2}$ and $\mathcal{M}_{k}^{o}=-2 k \partial_{x}$. By using

$$
\begin{equation*}
\phi(x)=\frac{e^{i k x}}{2}\left(u_{1}-i u_{2}\right)+\frac{e^{-i k x}}{2}\left(u_{1}+i u_{2}\right) \tag{11.28}
\end{equation*}
$$

and via simple computations, we obtain

$$
\mathcal{M} \phi=\cos (k x)\left(\mathcal{M}_{k}^{e} u_{1}+\mathcal{M}_{k}^{o} u_{2}\right)+\sin (k x)\left(-\mathcal{M}_{k}^{o} u_{1}+\mathcal{M}_{k}^{e} u_{2}\right),
$$

and

$$
J \phi=\cos (k x)\left(\partial_{x} u_{1}+k u_{2}\right)+\sin (k x)\left(\partial_{x} u_{2}-k u_{1}\right) .
$$

Define the operators

$$
J_{k}=\left(\begin{array}{cc}
\partial_{x} & k  \tag{11.29}\\
-k & \partial_{x}
\end{array}\right), \quad L_{k}=\left(\begin{array}{cc}
\mathcal{M}_{k}^{e}+c-f^{\prime}\left(u_{c, a}\right) & \mathcal{M}_{k}^{o} \\
-\mathcal{M}_{k}^{o} & \mathcal{M}_{k}^{e}+c-f^{\prime}\left(u_{c, a}\right)
\end{array}\right)
$$

Then $J_{k}, L_{k}$ are skew-adjoint and self-adjoint real operators and

$$
\mathcal{J} \mathcal{L}_{0} \phi=(\cos (k x), \sin (k x)) J_{k} L_{k}\binom{u_{1}}{u_{2}} .
$$

As always in the spectral analysis, $u_{1}$ and $u_{2}$, as well as operator $J_{k} L_{k}$ and quadratic forms $\left\langle L_{k} \cdot, \cdot\right\rangle$ and $\left\langle\cdot, J_{k} \cdot\right\rangle$, need to be complexified. By using operators $\mathcal{J}_{k}$ and $\mathcal{L}_{k}$ we can diagonalize $J_{k} L_{k}$ and $L_{k}$ blockwisely. In fact, let

$$
w_{1}=\frac{1}{2}\left(u_{1}-i u_{2}\right), \quad w_{2}=\frac{1}{2}\left(u_{1}+i u_{2}\right), \quad S\binom{u_{1}}{u_{2}}=\binom{w_{1}}{w_{2}} .
$$

One may compute using (11.28) and the definition of $\mathcal{J}_{k}$ and $\mathcal{L}_{k}$

$$
L_{k}=S^{-1}\left(\begin{array}{cc}
\mathcal{L}_{k} & 0  \tag{11.30}\\
0 & \mathcal{L}_{-k}
\end{array}\right) S, \quad J_{k} L_{k}=S^{-1}\left(\begin{array}{cc}
\mathcal{J}_{k} \mathcal{L}_{k} & 0 \\
0 & \mathcal{J}_{-k} \mathcal{L}_{-k}
\end{array}\right) S
$$

Moreover, $\mathcal{L}_{-k}$ and $\mathcal{J}_{-k} \mathcal{L}_{-k}$ are the complex conjugates of $L_{k}$ and $J_{k} L_{k}$ respectively, namely,

$$
\begin{equation*}
\mathcal{L}_{-k} w=\overline{\mathcal{L}_{k} \bar{w}}, \quad \mathcal{J}_{-k} \mathcal{L}_{-k} w=\overline{\mathcal{J}_{k} \mathcal{L}_{k} \bar{w}} \tag{11.31}
\end{equation*}
$$

From the above relations, we obtain

$$
n^{-}\left(L_{k}\right)=n^{-}\left(\mathcal{L}_{k}\right)+n^{-}\left(\mathcal{L}_{-k}\right)=2 n^{-}\left(\mathcal{L}_{k}\right)
$$

where $n^{-}\left(\mathcal{L}_{k}\right)$ is understood as the negative index of the complex Hermitian form $\left\langle\mathcal{L}_{k}, \cdot \cdot\right\rangle$. Moreover, (11.31) implies that $\lambda \in \sigma\left(\mathcal{J}_{k} \mathcal{L}_{k}\right)$ if and only if $\bar{\lambda} \in \sigma\left(\mathcal{J}_{-k} \mathcal{L}_{-k}\right)$, with $\operatorname{ker}\left(\bar{\lambda}-\mathcal{J}_{-k} \mathcal{L}_{-k}\right)^{n}$ consisting exactly of the complex conjugates of the functions in $\operatorname{ker}\left(\lambda-\mathcal{J}_{k} \mathcal{L}_{k}\right)^{n}$ for any $n>0$. Next, it is easy to see that $J_{k} L_{k}$, as well as $\mathcal{J}_{k} \mathcal{L}_{k}$, has compact resolvents and thus $\sigma\left(J_{k} L_{k}\right)$, as well as $\sigma\left(\mathcal{J}_{k} \mathcal{L}_{k}\right)$, consists of only discrete eigenvalues of finite algebraic multiplicity. Therefore, by Proposition 2.3, for any purely imaginary eigenvalue $i \mu \in \sigma\left(J_{k} L_{k}\right), L_{k}$ is non-degenerate on the finite dimensional eigenspace $E_{i \mu}$, and thus $n^{\leq 0}\left(\left.L_{k}\right|_{E_{i \mu}}\right)=n^{-}\left(\left.L_{k}\right|_{E_{i \mu}}\right)$. Let $\left(k_{r}, k_{c}, k_{i}^{-}, k_{0}^{-}\right)$be the indices defined in (2.10), (2.11), and (2.12) for $J_{k} L_{k}$, and $\left(\tilde{k}_{r}, \tilde{k}_{0}^{-}\right)$be the corresponding indices for the positive and zero eigenvalues of $\mathcal{J}_{k} \mathcal{L}_{k}$. Let $\tilde{k}_{c}$ be the sum of algebraic multiplicities of eigenvalues of $\mathcal{J}_{k} \mathcal{L}_{k}$ in the first and the fourth quadrants, $\tilde{k}_{i}^{-}$be the total number of negative dimensions of $\left\langle\mathcal{L}_{k}, \cdot,\right\rangle$ restricted to the subspaces of generalized eigenvectors of nonzero purely imaginary eigenvalues of $\mathcal{J}_{k} \mathcal{L}_{k}$. On the one hand, (11.30) and (11.31) imply

$$
k_{r}=2 \tilde{k}_{r}, \quad k_{c}=\tilde{k}_{c}, \quad k_{i}^{-}=\tilde{k}_{i}^{-}, \quad k_{0}^{-}=2 \tilde{k}_{0}^{-}
$$

On the other hand, Theorem 2.3 implies

$$
\begin{equation*}
k_{r}+2 k_{c}+2 k_{i}^{-}+k_{0}^{-}=2 n^{-}\left(\mathcal{L}_{k}\right) . \tag{11.32}
\end{equation*}
$$

Therefore, we obtain
Proposition 11.2. For any $k \in(0,1)$,

$$
\begin{equation*}
\tilde{k}_{r}+\tilde{k}_{c}+\tilde{k}_{i}^{-}+\tilde{k}_{0}^{-}=n^{-}\left(\mathcal{L}_{k}\right) \tag{11.33}
\end{equation*}
$$

The modulational instability occurs if $\tilde{k}_{r} \neq 0$ or $\tilde{k}_{c} \neq 0$.
Remark 11.1. Note that $\mathcal{J}_{k}$ is invertible for any $k \notin \mathbf{Z}$. With a more concrete form of $\mathcal{M}$, it is possible to determine $\tilde{k}_{0}^{-}$.

- Firstly, if $\operatorname{ker} \mathcal{L}_{0}$ is known (recall $\partial_{x} u_{c, a} \in \operatorname{ker} \mathcal{L}_{0}$ ), then one may study $\operatorname{ker} \mathcal{L}_{k}$, as well as $\tilde{k}_{0}^{-}$, for $0<|k| \ll 1$ through asymptotic analysis.
- If $\mathcal{M}=-\partial_{x x}$, then $\operatorname{ker} \mathcal{L}_{k}=\{0\}$ for any $k \in(0,1)$ (and thus for any $k \notin \mathbf{Z}$ ). In fact, in this case,

$$
\mathcal{L}_{0}=-\partial_{x x}+c-f^{\prime}\left(u_{c, a}\right), \quad v \in \operatorname{ker} \mathcal{L}_{k} \Longleftrightarrow e^{i k x} v \in \operatorname{ker} \mathcal{L}_{0}
$$

and $\operatorname{ker} \mathcal{L}_{0}=\operatorname{span}\left\{\partial_{x} u_{c, a}\right\}$. Suppose $\mathcal{L}_{k}$ has nontrivial kernel for some $k \in(0,1)$ and $0 \neq v \in \operatorname{ker} \mathcal{L}_{k}$. Denote $v_{0} \triangleq \partial_{x} u_{c, a}$, then the Wronskian of $v_{0}$ and $e^{i k x} v$ satisfies

$$
W(x)=e^{i k x}\left(v_{x} v_{0}-v v_{0 x}+i k v v_{0}\right)=\text { const } .
$$

Since $v$ and $v_{0}$ are $2 \pi$-periodic and $k \in(0,1)$, it must hold that

$$
\begin{equation*}
v_{x} v_{0}-v v_{0 x}+i k v v_{0}=0 \tag{11.34}
\end{equation*}
$$

We $\operatorname{claim} v(x) \neq 0$ for any $x \in \mathbf{S}^{1}$. In fact, if $v\left(x_{0}\right)=0$, then $v_{x}\left(x_{0}\right) \neq 0$ and (11.34) imply $v_{0}\left(x_{0}\right)=0$. The uniqueness of the solution to the ODE $\mathcal{L}_{0} u=0$ leads to the proportionality between $v_{0}$ and $e^{i k x} v$, a contradiction to $k \in(0,1)$ and the $2 \pi$-periodicity of $v(x)$. Now that $v(x) \neq 0$, (11.34) implies $\frac{v_{0}}{v}=C e^{-i k x}$, which is again a contradiction.

REMARK 11.2. The above index formula (11.33) was proved in [32] for the case when $\operatorname{ker} \mathcal{L}_{k}=\{0\}$, with additional assumptions to ensure that the generalized eigenfunctions of $\mathcal{J}_{k} \mathcal{L}_{k}$ form a complete basis of $L^{2}\left(\mathbf{S}^{1}\right)$ as assumed in the case of co-periodic perturbations. Proposition 11.2 is proved without such assumptions.

REmARK 11.3. We can also consider the case when the operator $\mathcal{M}$ is a smoothing operator, that is, $\|\mathcal{M}(\cdot)\|_{H^{r}} \sim\|\cdot\|_{L^{2}}$ for some $r>0$. One example is the Whitham equation which is a KDV type equation (11.2) with the symbol of $\mathcal{M}$ being $\sqrt{\frac{\tanh \xi}{\xi}}$ and thus $r=\frac{1}{2}$. In this case, if we assume that

$$
\begin{equation*}
-c-\left\|f^{\prime}\left(u_{c}\right)\right\|_{L^{\infty}\left(\mathbb{T}_{2 \pi}\right)} \geqslant \epsilon>0 \tag{11.35}
\end{equation*}
$$

then $\mathcal{L}_{0}$ and $\mathcal{L}_{k}$ are compact perturbations of the positive operator $-c+f^{\prime}\left(u_{c, a}\right)$ so that $n^{-}\left(-\mathcal{L}_{0}\right), n^{-}\left(-\mathcal{L}_{k}\right)<\infty$. Then the index formula

$$
\bar{k}_{r}+\bar{k}_{c}+k_{i}^{-}+k_{0}^{-}=n^{-}\left(-\mathcal{L}_{k}\right)
$$

is still true for the operator $\mathcal{J}_{k} \mathcal{L}_{k}, k \in(0,1)$. The assumption (11.35) can be verified ([36]) for small amplitude periodic traveling waves of Whitham equation with $f(u)=u^{2}$.

Under the conditions of orbital stability in Proposition 11.1, the spectra of the operator $J \mathcal{L}_{0}$ in $L^{2}\left(\mathbf{S}^{1}\right)$ lie on the imaginary axis and are all discrete. Moreover, the non-degeneracy of the matrix $D$ (defined by (11.21)) implies that the generalized kernel of $J \mathcal{L}_{0}$ is spanned by $\left\{\partial_{x} u_{c, a}, \partial_{c} u_{c, a}, U_{c, a}\right\}$. For $k \in(0,1)$ small, it is natural to study the spectra of $\mathcal{J}_{k} \mathcal{L}_{k}$ by the perturbation theory. Even though the results in Section 2.5 and Chapter 9 do not apply directly as $\mathcal{J}_{k}-J: X^{*} \rightarrow X$ is not bounded, the ideas there and the property that $\mathcal{J}_{k} \mathcal{L}_{k}$ has only isolated eigenvalues still yield the desired results. We start with the following lemma on the resolvent of $\mathcal{J}_{k} \mathcal{L}_{k}$.

Lemma 11.1. Assume that the symbol $\alpha(\xi)$ of $\mathcal{M}$ satisfies $a|\xi|^{m} \leq \alpha(\xi) \leq b|\xi|^{m}$, $a, b>0, m>0$, for large $\xi$ and

$$
\begin{equation*}
\lim _{\rho \rightarrow 0} \sup _{\xi \in \mathbf{Z}} \frac{|\alpha(\xi+\rho)-\alpha(\xi)|}{1+|\xi|^{m}} \rightarrow 0 \tag{11.36}
\end{equation*}
$$

then the resolvent $\left(\lambda-\mathcal{J}_{k} \mathcal{L}_{k}\right)^{-1}$ is continuous in $k \in[0,1]$.
Proof. Fix $k \in[0,1]$. From (11.26), one can compute

$$
\mathcal{J}_{k^{\prime}} \mathcal{L}_{k^{\prime}}-\mathcal{J}_{k} \mathcal{L}_{k}=\left(\partial_{x}+i k\right)\left(\mathcal{M}_{k^{\prime}}-\mathcal{M}_{k}\right)+i\left(k^{\prime}-k\right)\left(\mathcal{M}_{k^{\prime}}+c-f^{\prime}\left(u_{c, a}\right)\right) .
$$

On the one hand, there exists $a_{0} \neq 0$ such that $a_{0}+\left(\partial_{x}+i k\right) \mathcal{M}_{k}$ has a compact inverse on $X$. We obtain from (11.36)

$$
\begin{equation*}
\left|\left(a_{0}+\left(\partial_{x}+i k\right) \mathcal{M}_{k}\right)^{-1}\left(\mathcal{J}_{k^{\prime}} \mathcal{L}_{k^{\prime}}-\mathcal{J}_{k} \mathcal{L}_{k}\right)\right| \rightarrow 0 \text { as } k^{\prime} \rightarrow k \tag{11.37}
\end{equation*}
$$

On the other hand, (11.26) and $m>0$ imply that

$$
\begin{aligned}
& I+\left(a_{0}+\left(\partial_{x}+i k\right) \mathcal{M}_{k}\right)^{-1}\left(\lambda-\mathcal{J}_{k} \mathcal{L}_{k}\right) \\
= & \left(I+\left(\partial_{x}+i k\right) \mathcal{M}_{k}\right)^{-1}\left(\lambda+a_{0}-\left(\partial_{x}+i k\right)\left(c-f^{\prime}\left(u_{c, a}\right)\right)\right.
\end{aligned}
$$

is compact. Therefore, $A=\left(a_{0}+\left(\partial_{x}+i k\right) \mathcal{M}_{k}\right)^{-1}\left(\lambda-\mathcal{J}_{k} \mathcal{L}_{k}\right)$ is a Fredholm operator of index 0 . Suppose $\lambda \notin \sigma\left(\mathcal{J}_{k} \mathcal{L}_{k}\right)$, then $A$ is injective and thus $A^{-1}$ is bounded on
$X$. Along with (11.37), we obtain
$\left|\left(\lambda-\mathcal{J}_{k} \mathcal{L}_{k}\right)^{-1}\left(\mathcal{J}_{k^{\prime}} \mathcal{L}_{k^{\prime}}-\mathcal{J}_{k} \mathcal{L}_{k}\right)\right|=\left|A^{-1}\left(a_{0}+\left(\partial_{x}+i k\right) \mathcal{M}_{k}\right)^{-1}\left(\mathcal{J}_{k^{\prime}} \mathcal{L}_{k^{\prime}}-\mathcal{J}_{k} \mathcal{L}_{k}\right)\right| \rightarrow 0$ as $k^{\prime} \rightarrow k$. From

$$
\lambda-\mathcal{J}_{k^{\prime}} \mathcal{L}_{k^{\prime}}=\left(\lambda-\mathcal{J}_{k} \mathcal{L}_{k}\right)\left(I-\left(\lambda-\mathcal{J}_{k} \mathcal{L}_{k}\right)^{-1}\left(\mathcal{J}_{k^{\prime}} \mathcal{L}_{k^{\prime}}-\mathcal{J}_{k} \mathcal{L}_{k}\right)\right)
$$

we obtain the continuity of the resolvent $\left(\lambda-\mathcal{J}_{k} \mathcal{L}_{k}\right)^{-1}$ in $k \in[0,1]$.
REmARK 11.4. The assumption (11.36) is clearly satisfied if $\alpha(\xi) \in C^{1}(\mathbf{R})$ and

$$
\limsup _{|\xi| \rightarrow \infty} \frac{\alpha^{\prime}(\xi)}{|\xi|^{m}}<\infty
$$

Next we show that when $k$ is small enough, the unstable modes of $\mathcal{J}_{k} \mathcal{L}_{k}$ can only bifurcate from the zero eigenvalue of $J \mathcal{L}_{0}$.

Proposition 11.3. Suppose $\operatorname{ker} \mathcal{L}_{0}=\operatorname{span}\left\{\partial_{x} u_{c, a}\right\}, D$ is nonsingular, $n^{-}\left(\mathcal{L}_{0}\right)=$ $n^{-}(D)$ and (11.36) holds. Then for any $\delta>0$, there exists $\varepsilon_{0}>0$ such that if $|k|<\varepsilon_{0}$, then $\sigma\left(\mathcal{J}_{k} \mathcal{L}_{k}\right) \cap\{|z| \geq \delta\} \subset i \mathbf{R}$.

Proof. Since 0 is an isolated spectral point of $J \mathcal{L}_{0}$, there exists $\delta_{0}>0$ such that $\lambda \notin \sigma\left(J \mathcal{L}_{0}\right)$ as long as $0<|\lambda| \leq \delta_{0}$. Without loss of generality, assume $0<\delta<\delta_{0}$. Lemma 11.1 implies $\lambda \notin \sigma\left(\mathcal{J}_{k} \mathcal{L}_{k}\right)$ for $0<|k| \ll 1$. Let

$$
P(k)=\frac{1}{2 \pi i} \oint_{|\lambda|=\delta}\left(\lambda-\mathcal{J}_{k} \mathcal{L}_{k}\right)^{-1} d \lambda, \quad Z_{k}=P(k) X, \quad Y_{k}=(I-P(k)) X
$$

The standard spectral theory implies that $P(k)$ is continuous in $k, Y_{k}$ and $Z_{k}$ are invariant under $\mathcal{J}_{k} \mathcal{L}_{k}$, and

$$
|\lambda|<\delta, \forall \lambda \in \sigma\left(\left.\mathcal{J}_{k} \mathcal{L}_{k}\right|_{Z_{k}}\right) \text { and }|\lambda|>\delta, \forall \lambda \in \sigma\left(\left.\mathcal{J}_{k} \mathcal{L}_{k}\right|_{Y_{k}}\right) .
$$

For $k=0$, our assumptions imply that $Z_{0}=\operatorname{span}\left\{\partial_{x} u_{c, a}, \partial_{c} u_{c, a}, U_{c, a}\right\}$. Therefore, $Z_{k}$ close to $Z_{0}$ is a 3-dim invariant subspace of $\mathcal{J}_{k} \mathcal{L}_{k}$ with small eigenvalues containing ker $\mathcal{L}_{k}$. Moreover, the assumption

$$
n^{-}\left(\mathcal{L}_{0}\right)=n^{-}(D)=n^{-}\left(\left.\mathcal{L}_{0}\right|_{Z_{0}}\right)
$$

and the $\mathcal{L}_{0}$-orthogonality between $Z_{0}$ and $Y_{0}$ imply that $\mathcal{L}_{0}$ is uniformly positive definite on $Y_{0}$. As $\mathcal{L}_{k}: X=H^{\frac{m}{2}} \rightarrow X^{*}=H^{-\frac{m}{2}}$ is continuous in $k$, there exists $\alpha>0$ such that $\left\langle\mathcal{L}_{k} u, u\right\rangle>\alpha\|u\|^{2}$ for all $u \in Y_{k}$. Clearly, $\left.\mathcal{J}_{k} \mathcal{L}_{k}\right|_{Y_{k}}$ is skewadjoint with respect to the equivalent inner product given by $\left\langle\mathcal{L}_{k} \cdot, \cdot\right\rangle$ on $Y_{k}$, therefore $\sigma\left(\left.\mathcal{J}_{k} \mathcal{L}_{k}\right|_{Y_{k}}\right) \subset i \mathbf{R}$ and the proposition follows.

Since $\operatorname{dim} \operatorname{ker}\left(J \mathcal{L}_{0}\right)=3$, the perturbation of zero eigenvalue of $J \mathcal{L}_{0}$ for $\mathcal{J}_{k} \mathcal{L}_{k}$ $(0<k \ll 1)$ can be reduced to the eigenvalue perturbation of a 3 by 3 matrix. This had been studied extensively in the literature and instability conditions were obtained for various dispersive models. See the survey [12] and the references therein.

Recently, it was proved in [36] that linear modulational instability of the traveling wave $u_{c}(x-c t)$ also implies the nonlinear instability for both multi-periodic and localized perturbations. The semigroup estimates of $e^{t J \mathcal{L}_{0}}$ play an important role on this proof of nonlinear instability. We sketch these estimates below, as an example of the application of Theorem 2.2 on the exponential trichotomy of linear Hamiltonian PDE. First, if $u_{c}$ is linearly modulationally unstable, then there exists a rational $k_{0}=\frac{p}{q} \in(0,1)$ such that $\mathcal{J}_{k_{0}} \mathcal{L}_{k_{0}}$ has an unstable eigenvalue. By the
definition of $\mathcal{J}_{k_{0}} \mathcal{L}_{k_{0}}$, this implies that the operator $J \mathcal{L}_{0}$ has an unstable eigenvalue on the $2 \pi q$ periodic space $L^{2}\left(\mathbf{S}_{2 \pi q}^{1}\right)$ with an eigenfunction of the form $e^{i k_{0} x} u(x)$ $\left(u \in L^{2}\left(\mathbf{S}^{1}\right)\right)$. The exponential trichotomy of the semigroup $e^{t J \mathcal{L}_{0}}$ on the space $H^{s}\left(\mathbf{S}_{2 \pi q}^{1}\right)\left(s \geq \frac{m}{2}\right)$ follows directly by Theorem 2.2. This is used in [36] to prove nonlinear orbital instability of $u_{c}$ for $2 \pi q$ periodic perturbations or even to construct stable and unstable manifolds. To prove nonlinear instability for localized perturbations, we study the semigroup $e^{t J \mathcal{L}_{0}}$ on the space $H^{s}(\mathbf{R})\left(s \geq \frac{m}{2}\right)$. The operator $\mathcal{L}_{0}$ might have negative continuous spectrum in $H^{s}(\mathbf{R})$. For example, when $\mathcal{M}=-\partial_{x}^{2}$, the spectrum of $\mathcal{L}_{0}=-\partial_{x}^{2}+V(x)$ with periodic $V(x)$ is well studied in the literature and is known to have bands of continuous spectrum. So Theorem 2.2 does not apply. However, we have the following upper bound estimate of $e^{t J \mathcal{L}_{0}}$ on $H^{s}(\mathbf{R})$, which suffices to prove nonlinear localized instability.

Lemma 11.2. Assume (11.36). Let $\lambda_{0} \geq 0$ be such that

$$
\operatorname{Re} \lambda \leq \lambda_{0}, \quad \forall \xi \in[0,1], \lambda \in \sigma\left(\mathcal{J}_{\xi} \mathcal{L}_{\xi}\right)
$$

For every $s \geq \frac{m}{2}$, there exist $C(s)>0$ such that

$$
\begin{align*}
& \left\|e^{t \mathcal{J}_{\xi} \mathcal{L}_{\xi}} v(x)\right\|_{H^{s}\left(\mathbf{S}^{1}\right)} \leq C(s)\left(1+t^{2 n^{-}\left(\mathcal{L}_{\xi}\right)+1}\right) e^{\lambda_{0} t}\|v(x)\|_{H^{s}\left(\mathbf{S}^{1}\right)}  \tag{11.38}\\
& \left\|e^{t J \mathcal{L}_{0}} u(x)\right\|_{H^{s}(\mathbf{R})} \leqslant C(s)\left(1+t^{2 n^{-}\left(\mathcal{L}_{\xi}\right)+1}\right) e^{\lambda_{0} t}\|u(x)\|_{H^{s}(\mathbf{R})} \tag{11.39}
\end{align*}
$$

for any $\xi \in[0,1], v \in H^{s}\left(\mathcal{S}^{1}\right)$, and $u \in H^{s}(\mathbf{R})$.
Proof. It suffices to prove the lemma for $s=\frac{m}{2}$. The estimates for general $s \geq \frac{m}{2}$ can be obtained by applying $\mathcal{J}_{\xi} \mathcal{L}_{\xi}$ and $J \mathcal{L}_{0}$ repeatedly to the estimates for $s=\frac{m}{2}$ (and interpolation for the case when $\frac{2 s}{m}$ is not an integer). We start with the first estimate in the $2 \pi$-periodic case. Due to the compactness of $[0,1]$, it suffices to prove that for any $\xi_{0} \in[0,1]$, there exist $C, \epsilon>0$ and an integer $K \geq 0$ such that (11.38) holds for $\xi \in\left(\xi_{0}-\epsilon, \xi_{0}+\epsilon\right)$. We first note that each $\lambda \in \sigma\left(\mathcal{J}_{\xi_{0}} \mathcal{L}_{\xi_{0}}\right)$ is an isolated eigenvalue with finite algebraic multiplicity and $\mathcal{L}_{\xi_{0}}$ is non-degenerate on $E_{\lambda} /\left(E_{\lambda} \cap \operatorname{ker} \mathcal{L}_{\xi_{0}}\right)$. Let

$$
\Lambda=\left\{\lambda \in \sigma\left(\mathcal{J}_{\xi_{0}} \mathcal{L}_{\xi_{0}}\right) \mid \exists \delta>0 \text { s.t. }\left\langle\mathcal{L}_{\xi_{0}} v, v\right\rangle \geq \delta\|v\|^{2}\right\}
$$

Due to Proposition 11.2, $\sigma\left(\mathcal{J}_{\xi_{0}} \mathcal{L}_{\xi_{0}}\right) \backslash \Lambda$ is finite and

$$
n=\Sigma_{\lambda \in \sigma\left(\mathcal{J}_{\xi_{0}} \mathcal{L}_{\xi_{0}}\right) \backslash \Lambda} \operatorname{dim} E_{\lambda}<\infty
$$

Moreover, there exists $\varepsilon>0$ such that

$$
\Omega \cap \Lambda=\emptyset, \quad \text { where } \Omega=\cup_{\lambda \in \sigma\left(\mathcal{J}_{\xi_{0}} \mathcal{L}_{\xi_{0}}\right) \backslash \Lambda}\{z \mid\|z-\lambda\|<\varepsilon\} \subset \mathbf{C} .
$$

From Lemma 11.1, there exists $\epsilon>0$ such that $\partial \Omega \cap \sigma\left(\mathcal{J}_{\xi} \mathcal{L}_{\xi}\right)=\emptyset$ for any $\xi \in$ $\left[\xi_{0}-\epsilon, \xi_{0}+\epsilon\right]$. For such $\xi$, let

$$
P(\xi)=\frac{1}{2 \pi i} \oint_{\partial \Omega}\left(\lambda-\mathcal{J}_{k} \mathcal{L}_{k}\right)^{-1} d \lambda, \quad Z_{\xi}=P(\xi) X, \quad Y_{\xi}=(I-P(\xi)) X
$$

which are continuous in $\xi$ and invariant under $e^{t \mathcal{J}_{\xi} \mathcal{L}_{\xi}}$. Therefore, $\operatorname{dim} Z_{\xi}=n$ and the continuity of $\mathcal{L}_{\xi}$ in $\xi$ implies that there exists $\delta>0$ such that

$$
\delta^{-2}\|v\|^{2} \geq\left\langle\mathcal{L}_{\xi} v, v\right\rangle \geq \delta^{2}\|v\|^{2}, \quad \forall v \in Y_{\xi},\left|\xi-\xi_{0}\right| \leq \epsilon
$$

Moreover, according to Proposition 2.2, for any $\lambda \in \Omega \cap \sigma\left(\mathcal{J}_{\xi} \mathcal{L}_{\xi}\right)$, the dimension of its eigenspace

$$
E_{\lambda}\left(\mathcal{J}_{\xi} \mathcal{L}_{\xi}\right)=\operatorname{ker}\left(\lambda-\mathcal{J}_{\xi} \mathcal{L}_{\xi}\right)^{2\left(1+n^{-}\left(\mathcal{L}_{\xi}\right)\right)}
$$

namely, the maximal dimension of Jordan blocks of $\mathcal{J}_{\xi} \mathcal{L}_{\xi}$ on $Y_{\xi}$ is no more than $2\left(1+n^{-}\left(\mathcal{L}_{\xi}\right)\right)$. So for any $\xi \in\left[\xi_{0}-\epsilon, \xi_{0}+\epsilon\right]$, there exists a generic constant $C>0$ independent of $\xi$, such that

$$
\begin{aligned}
& \left\|e^{t \mathcal{J}_{\xi} \mathcal{L}_{\xi}} v\right\| \leq\left\|e^{t \mathcal{J}_{\xi} \mathcal{L}_{\xi}} P(\xi) v\right\|+\left\|e^{t \mathcal{J}_{\xi} \mathcal{L}_{\xi}}(I-P(\xi)) v\right\| \\
\leq & C\left(\left(1+t^{2 n^{-}\left(\mathcal{L}_{\xi}\right)+1}\right) e^{\lambda_{0} t}\|P(\xi) v\|+\left\langle\mathcal{L}_{\xi} e^{t \mathcal{J}_{\xi} \mathcal{L}_{\xi}}(I-P(\xi)) v, e^{t \mathcal{J}_{\xi} \mathcal{L}_{\xi}}(I-P(\xi)) v\right\rangle^{\frac{1}{2}}\right) \\
\leq & C\left(\left(1+t^{2 n^{-}\left(\mathcal{L}_{\xi}\right)+1}\right) e^{\lambda_{0} t}\|P(\xi) v\|+\left\langle\mathcal{L}_{\xi}(I-P(\xi)) v,(I-P(\xi)) v\right\rangle^{\frac{1}{2}}\right) \\
\leq & C\left(1+t^{2 n^{-}\left(\mathcal{L}_{\xi}\right)+1}\right) e^{\lambda_{0} t}\|v\| .
\end{aligned}
$$

Along with the compactness of $[0,1]$, it implies (11.38).
To prove (11.39), we first write, for any $u \in H^{s}(\mathbf{R})$,

$$
u(x)=\int_{0}^{1} e^{i \xi x} u_{\xi}(x) d \xi, \text { where } u_{\xi}(x)=\Sigma_{n \in \mathbf{Z}} e^{i n x} \hat{u}(n+\xi) \in H^{s}\left(\mathbf{S}^{1}\right)
$$

and $\hat{u}$ is the Fourier transform of $u$. Clearly, there exists $C>0$ such that

$$
\begin{equation*}
\frac{1}{C}\|u\|_{H^{s}(\mathbb{R})}^{2} \leq \int_{0}^{1}\left\|u_{\xi}(x)\right\|_{H^{s}\left(\mathbf{S}^{1}\right)}^{2} d \xi \leq C\|u\|_{H^{s}(\mathbb{R})}^{2} \tag{11.40}
\end{equation*}
$$

Note

$$
e^{t J \mathcal{L}_{0}} u(x)=\int_{0}^{1} e^{i \xi x} e^{t \mathcal{J}_{\xi} \mathcal{L}_{\xi}} u_{\xi}(x) d \xi
$$

and thus

$$
\begin{equation*}
\left\|e^{t J \mathcal{L}_{0}} u(x)\right\|_{H^{s}(\mathbf{R})}^{2} \approx \int_{0}^{1}\left\|e^{t \mathcal{J}_{\xi} \mathcal{L}_{\xi}} u_{\xi}(x)\right\|_{H^{s}\left(\mathbf{S}^{1}\right)}^{2} d \xi \tag{11.41}
\end{equation*}
$$

Along with (11.38), it immediately implies (11.39).
REMARK 11.5. The semigroup estimates of the types (11.38) and (11.39) can also be obtained for $s=-1$, that is, in the negative Sobolev space $H^{-1}\left(\mathbf{S}_{2 \pi q}^{1}\right)$ and $H^{-1}(\mathbf{R})$ for $e^{t J \mathcal{L}_{0}}$ (see [36]). Such semigroup estimates were used in [36] to prove nonlinear modulational instability by a bootstrap argument.
11.4. The spectral problem $L u=\lambda u^{\prime}$

In this section, we consider the eigenvalue problem of the form

$$
\begin{equation*}
L u=\lambda u^{\prime} \tag{11.42}
\end{equation*}
$$

where the symmetric operator $L$ is of the form of $\mathcal{L}_{0}$ in Section 11.1. As an example, consider the stability of solitary waves of generalized Bullough-Dodd equation ([69])

$$
\begin{equation*}
u_{t x}=a u-f(u) \tag{11.43}
\end{equation*}
$$

where $a>0$ and $f$ is a smooth function of $u$ satisfying

$$
\begin{equation*}
f(u)=O\left(u^{2}\right), f^{\prime}(u)=O(u) \text { for small } u \tag{11.44}
\end{equation*}
$$

The traveling wave $u_{c}(x+c t)$ satisfies the ODE

$$
-c u_{c}^{\prime \prime}+a u_{c}-f\left(u_{c}\right)=0
$$

Then the linearized equation in the traveling frame $(x+c t, t)$ takes the form

$$
\begin{equation*}
u_{t x}=-c u_{x x}+a u-f^{\prime}\left(u_{c}\right) u \tag{11.45}
\end{equation*}
$$

Thus the eigenvalue problem takes the form (11.42) with

$$
\begin{equation*}
L=-c \frac{d^{2}}{d x^{2}}+a-f^{\prime}\left(u_{c}\right) \tag{11.46}
\end{equation*}
$$

We consider the general problem (11.42) with $L$ of the form $L=\mathcal{M}+V(x)$. We assume that: i) $M$ is a Fourier multiplier operator with the symbol $\alpha(\xi)$ satisfying

$$
\begin{equation*}
\alpha(\xi) \geq 0 \text { and } \alpha(\xi) \approx|\xi|^{2 s}(s>0), \text { when }|\xi| \text { is large }, \tag{11.47}
\end{equation*}
$$

and ii) the real potential $V(x)$ satisfies

$$
\begin{equation*}
V(x) \rightarrow \delta_{0}>0 \text { when }|x| \rightarrow \infty \tag{11.48}
\end{equation*}
$$

Let $X=H^{s}(R)(s>0)$. Then the assumption (H2) is satisfied for $L$ on $X$. Namely, $L: X \rightarrow X^{*}$ is bounded and symmetric, and there exists a decomposition of $X$

$$
X=X_{-} \oplus \operatorname{ker} L \oplus X_{+}, \quad n^{-}(L) \triangleq \operatorname{dim} X_{-}<\infty
$$

satisfying $\left.L\right|_{X_{-}}<0$ and $\left.L\right|_{X_{+}} \geq \delta>0$.
Define $J=\partial_{x}^{-1}$. Now we check that $J: X^{*} \rightarrow X$ is densely defined and $J^{*}=-J$. On $X=H^{s}(R)$ with $s>0$, the operator $\partial_{x}: X \rightarrow X^{*}$ is densely defined and satisfies $\left(\partial_{x}\right)^{*}=-\partial_{x}$. Since ker $\partial_{x}=\{0\}$,

$$
\overline{R\left(\partial_{x}\right)}=\left(\operatorname{ker}\left(\partial_{x}^{*}\right)\right)^{\perp}=\left(\operatorname{ker}\left(-\partial_{x}\right)\right)^{\perp}=X^{*},
$$

so $D\left(\partial_{x}^{-1}\right)=R\left(\partial_{x}\right)$ is dense in $X^{*}$ and $J=\partial_{x}^{-1}: X^{*} \rightarrow X$ satisfies $J^{*}=-J$.
So the eigenvalue problem $L u=\lambda u^{\prime}$ can be equivalently written in the Hamiltonian form $J L u=\lambda u$, where $(J, L, X)$ satisfies the assumptions (H1)-(H3). Let $\operatorname{ker} L=\operatorname{span}\left\{\psi_{1}, \cdots, \psi_{l}\right\}$ and

$$
\operatorname{span}\left\{\psi_{1}^{\prime}, \cdots, \psi_{l}^{\prime}\right\} \cap R(L)=\operatorname{span}\left\{g_{1}, \cdots, g_{m}\right\}, m \leq l .
$$

Define the $m$ by $m$ matrix

$$
D=\left(\left\langle L^{-1} g_{i}, g_{j}\right\rangle\right), 1 \leq i, j \leq m
$$

By Theorem 2.3 and Proposition 2.7, we get the following theorem.
Theorem 11.3. Assume (11.47) and (11.48). Then

$$
k_{r}+2 k_{c}+2 k_{i}^{\leq 0}+k_{0}^{\leq 0}=n^{-}(L),
$$

where $k_{r}, k_{c}, k_{i}^{\leq 0}, k_{0}^{\leq 0}$ are the indexes for the eigenvalues of $\partial_{x}^{-1} L$, as defined in Section 2.4. In addition, we have $k_{0}^{\leq 0} \geq n^{\leq 0}(D)$, where $n^{\leq 0}(D)$ is the number of nonpositive eigenvalues of $D$. If $D$ is nonsingular, then $k_{0}^{\leq 0}=n^{-}(D)$, i.e., the number of negative eigenvalues of $D$.

For many applications, particularly the generalized Bullough-Dodd equation where $\mathcal{M}=-c \partial_{x}^{2}(c>0), L$ has at most one dimensional kernel and negative eigenspace. In this case, we get a more explicit instability criterion.

Corollary 11.4. i) Assume $n^{-}(L)=1$ and $\operatorname{ker} L=\left\{\psi_{0}\right\}$. Then there is a positive eigenvalue of $\partial_{x}^{-1} L$ when $\left\langle L^{-1} \psi_{0}^{\prime}, \psi_{0}^{\prime}\right\rangle>0$.
ii) Assume $n^{-}(L) \leq 1$ and there exists $0 \neq \psi_{0} \in \operatorname{ker} L$ such that $\left\langle L^{-1} \psi_{0}^{\prime}, \psi_{0}^{\prime}\right\rangle \leq$ 0 , then $\partial_{x}^{-1} L$ has no unstable eigenvalues.

Remark 11.6. The above Corollary was obtained in [69] under some additional assumptions. In [69], Corollary 11.4 i) was proved under the following two assumptions:

C1) $\left(f_{0}, g_{0}\right) \neq 0$, where $f_{0}$ is the eigenfunction of $L$ with the negative eigenvalue and $g_{0}^{\prime} \in \operatorname{ker} L$.

C2) For any $\lambda \in \mathbf{R}$,

$$
\left\|P_{+}\left(L-\lambda \partial_{x}\right)^{-1} P_{+} v\right\|_{H^{1}} \leq C(\lambda)\|v\|_{L^{2}}
$$

where $P_{+}$is the projection to the positive space of $L$ and $C(\lambda)$ is bounded on compact sets.
The proof in [69] is by constructing Evans-like functions. Corollary 11.4 ii) was proved in [69] under the following additional assumptions:

D1) $\operatorname{ker} L=\left\{\psi_{0}\right\}$ and $\left\langle L^{-1} \psi_{0}^{\prime}, \psi_{0}^{\prime}\right\rangle<0$;
D2) For any $\lambda \notin i \mathbf{R}$, the operator $L-\lambda \partial_{x}$ has zero index and the equation $\left(L-\lambda \partial_{x}\right) f=g$ satisfies certain Fredholm alternative properties (see (12)(13)(14) in [69]);

D3) The symbol $\alpha(\xi)$ of the leading order part $\mathcal{M}$ of $L$ satisfies

$$
\alpha(\xi) \approx|\xi|^{2 s}\left(s>\frac{1}{2}\right), \text { when }|\xi| \text { is large. }
$$

The proof in [69] is by Lyapunov-Schmidt reduction arguments and the index theorem in [44].

For the Bullough-Dodd equation (11.43), ker $L=\left\{u_{c, x}\right\}$ where $L$ is defined by (11.46). Since the momentum of the problem is $\frac{1}{2} \int\left(u_{c}^{\prime}\right)^{2} d x$, by similar computation as in (11.12), it was shown in [69] that

$$
\left\langle L^{-1} u_{c}^{\prime \prime}, u_{c}^{\prime \prime}\right\rangle=-\frac{1}{2} \partial_{c} \int\left(u_{c}^{\prime}\right)^{2} d x=-\frac{1}{2} \partial_{c}\left[c^{-\frac{1}{2}} \int\left(u_{1}^{\prime}\right)^{2} d x\right]>0
$$

where $u_{c}=u_{1}(x / \sqrt{c})$ and $-u_{1}^{\prime \prime}+a u_{1}-f\left(u_{1}\right)=0$. So we get the following
THEOREM 11.4. Assume $f(u)$ is a smooth function satisfying (11.44) and the traveling wave solution $u_{c}(x-c t)$ to (11.43) exists with $c>0$ and $u_{c}(x) \rightarrow 0$ as $|x| \rightarrow \infty$, then $u_{c}$ is linearly unstable.

In [69], the above Theorem was proved for smooth and convex function $f$. Their additional convexity assumption on $f$ was used to verify the condition C1) in Remark 11.6.

Besides the above linear instability result, Theorem 2.2 can be applied to give the exponential trichotomy for the linearized equation (11.45). This will be useful for the construction of invariant manifolds of (11.43) near the unstable traveling wave orbit.

### 11.5. Stability of steady flows of 2 D Euler equation

We consider the 2D Euler equations

$$
\begin{gather*}
\partial_{t} u+(u \cdot \nabla u)+\nabla p=0  \tag{11.49}\\
\nabla \cdot u=0 \tag{11.50}
\end{gather*}
$$

in a bounded domain $\Omega \subset \mathbf{R}^{2}$ with smooth boundary $\partial \Omega$ composed of a finite number of connected components $\Gamma_{i}$. The boundary condition is

$$
u \cdot n=0 \quad \text { on } \quad \partial \Omega,
$$

For simplicity, first we consider $\Omega$ to be simply connected and $\partial \Omega=\Gamma$. The vorticity form of (11.49)-(11.50) is given by

$$
\begin{equation*}
\partial_{t} \omega-\psi_{y} \partial_{x} \omega+\psi_{x} \partial_{y} \omega=0 \tag{11.51}
\end{equation*}
$$

where $\psi$ is the stream function, then $\omega \equiv-\Delta \psi \equiv-\left(\partial_{x}^{2}+\partial_{y}^{2}\right) \psi$ is the vorticity and $u=\nabla^{\perp} \psi=\left(\psi_{y},-\psi_{x}\right)$ is the velocity. The boundary condition associated with (11.51) is given by $\psi=0$ on $\partial \Omega$. A stationary solution of (11.51) is given by a stream function $\psi_{0}$ satisfying

$$
\begin{equation*}
-\psi_{0_{y}} \partial_{x} \omega_{0}+\psi_{0_{x}} \partial_{y} \omega_{0}=0 \tag{11.52}
\end{equation*}
$$

here $\omega_{0} \equiv-\Delta \psi_{0}$ and $u_{0}=\nabla^{\perp} \psi_{0}$ are the associated vorticity and velocity. Suppose $\psi_{0}$ satisfy the following elliptic equation

$$
-\Delta \psi_{0}=g\left(\psi_{0}\right)
$$

with boundary condition $\psi_{0}=0$ on $\partial \Omega$, where $g$ is some differentiable function. Then $\omega_{0} \equiv-\Delta \psi_{0}=g\left(\psi_{0}\right)$ is a steady solution of equation (11.51). The linearized equation near $\omega_{0}$ is

$$
\begin{equation*}
\partial_{t} \omega-\psi_{0_{y}} \partial_{x} \omega+\psi_{0_{x}} \partial_{y} \omega=\psi_{y} \partial_{x} \omega_{0}-\psi_{x} \partial_{y} \omega_{0} \tag{11.53}
\end{equation*}
$$

with $\omega=-\Delta \psi$ and the boundary condition $\left.\psi\right|_{\partial \Omega}=0$. The above equation can be written as

$$
\begin{equation*}
\partial_{t} \omega-u_{0} \cdot \nabla \omega+g^{\prime}\left(\psi_{0}\right) u_{0} \cdot \nabla \psi=0 \tag{11.54}
\end{equation*}
$$

Below we consider the case when $g^{\prime}>0$ which appeared in many interesting cases such as mean field equations (e.g. [14] [15]). Then (11.54) has the following Hamiltonian structure

$$
\begin{equation*}
\partial_{t} \omega=J L \omega, \text { where } J=g^{\prime}\left(\psi_{0}\right) u_{0} \cdot \nabla, \quad L=\frac{1}{g^{\prime}\left(\psi_{0}\right)}-(-\Delta)^{-1} \tag{11.55}
\end{equation*}
$$

We take the energy space of the linearized Euler (11.55) as the weighted space

$$
X=\left\{\omega \mid\|\omega\|_{X}<\infty\right\}, \text { where }\|\omega\|_{X}=\left(\iint_{\Omega} \frac{|\omega|^{2}}{g^{\prime}\left(\psi_{0}\right)} d x d y\right)^{\frac{1}{2}}
$$

If $g^{\prime}$ has a positive lower bound, $X$ is equivalent to $L^{2}(\Omega)$. In general, $\omega \in X$ implies $\omega \in L^{2}$ and $\nabla \psi \in L^{2}$. Therefore, $\langle L \cdot, \cdot\rangle$ defines a bounded symmetric quadratic form on $X$ and $L: X \rightarrow X^{*}$ is a bounded symmetric operator. Moreover, it is easy to see that

$$
S: L^{2} \rightarrow X, \quad S \omega=g^{\prime}\left(\psi_{0}\right)^{\frac{1}{2}} \omega
$$

defines an isometry. As $f\left(\psi_{0}\right) \cdot$ and $u_{0} \cdot \nabla$ are commutative for any $f$, we have

$$
\tilde{J} \triangleq S^{-1} J\left(S^{*}\right)^{-1}=u_{0} \cdot \nabla:\left(L^{2}\right)^{*} \rightarrow L^{2}
$$

is anti-self-dual due to $\nabla \cdot u_{0}=0$, from which we obtain $J^{*}=-J$ and thus (H1) is satisfied by $J$ and $X$. Moreover, since $\frac{1}{g^{\prime}\left(\psi_{0}\right)} \cdot: X \rightarrow X^{*}$ is an isomorphism and
$(-\Delta)^{-1}$ is compact, we have $\operatorname{dim} \operatorname{ker} L<\infty$ and thus (H3) is satisfied. Note that the closed subspace ker $J \subset X^{*}$ is infinite dimensional since

$$
\operatorname{ker} J \supset\left\{h\left(\psi_{0}\right), h \in C^{1}\right\}
$$

Let $\tilde{P}:\left(L^{2}\right)^{*} \rightarrow \operatorname{ker} \tilde{J}$ be the orthogonal projection and define

$$
P=\left(S^{*}\right)^{-1} \tilde{P} S^{*}: X^{*} \rightarrow \operatorname{ker} J
$$

Clearly, $P$ is a bounded linear operator on $X^{*}$ and it defines a projection on $X^{*}$, but orthogonal in the $L^{2}$ sense. In fact, due to the commutativity between $f\left(\psi_{0}\right)$. and $u_{0} \cdot \nabla$ for any $f$, operators $P$ and $\tilde{P}$ take the same form shown in ([49])

$$
\left.P \phi\right|_{\gamma_{i}(c)}=\frac{\oint_{\gamma_{i}(c)} \frac{\phi(x, y)}{\left|\nabla \psi_{0}\right|} d l}{\oint_{\gamma_{i}(c)} \frac{1}{\left|\nabla \psi_{0}\right|} d l},
$$

where $c$ is in the range of $\psi_{0}$ and $\gamma_{i}(c)$ is a branch of $\left\{\psi_{0}=c\right\}$. As in ([49]), define operator $A$ : $H_{0}^{1} \cap H^{2}(\Omega) \rightarrow L^{2}(\Omega)$ by

$$
A \phi=-\Delta \phi-g^{\prime}\left(\psi_{0}\right) \phi+g^{\prime}\left(\psi_{0}\right) P \phi
$$

We also denote the operator

$$
A_{0}=-\Delta-g^{\prime}\left(\psi_{0}\right): H_{0}^{1} \cap H^{2}(\Omega) \rightarrow L^{2}(\Omega)
$$

Clearly, $A, A_{0}$ are self-adjoint with compact resolvents and thus with only discrete spectra. The next lemma studies the spectral information of $L$ on the weighted space $X$.

Recall that for any subspace $Y \in X,\langle L \cdot, \cdot\rangle$ also defines a bounded symmetric quadratic form on the quotient space $Y /(Y \cap \operatorname{ker} L)$.

Lemma 11.3. i) The assumption (H2) is satisfied by $\langle L \cdot, \cdot\rangle$ on $X$, with $n^{-}(L)=$ $n^{-}\left(A_{0}\right)$ and $\operatorname{dim} \operatorname{ker} L=\operatorname{dim} \operatorname{ker} A_{0}$.
ii) The quadratic form $\langle L \cdot, \cdot\rangle$ is non-degenerate on $\overline{R(J)} /(\overline{R(J)} \cap \operatorname{ker} L)$ if and only if $\operatorname{ker} A \subset \operatorname{ker} A_{0}$. Moreover,

$$
\begin{equation*}
n^{-}\left(\left.L\right|_{\overline{R(J)} /(\overline{R(J)} \cap \operatorname{ker} L)}\right)=n^{-}\left(\left.L\right|_{\overline{R(J)}}\right)=n^{-}(A) \tag{11.56}
\end{equation*}
$$

Proof. i) For any $\omega \in X$, we have

$$
\begin{align*}
\langle L \omega, \omega\rangle & =\iint_{\Omega}\left\{\frac{\omega^{2}}{g^{\prime}\left(\psi_{0}\right)}-|\nabla \psi|^{2}\right\} d x d y  \tag{11.57}\\
& =\iint_{\Omega}\left\{\frac{\omega^{2}}{g^{\prime}\left(\psi_{0}\right)}-2 \psi \omega+|\nabla \psi|^{2}\right\} d x d y \\
& =\iint_{\Omega}\left\{\left(\frac{\omega}{\sqrt{g^{\prime}\left(\psi_{0}\right)}}-\psi \sqrt{g^{\prime}\left(\psi_{0}\right)}\right)^{2}-g^{\prime}\left(\psi_{0}\right) \psi^{2}+|\nabla \psi|^{2}\right\} d x d y \\
& \geq \iint_{\Omega}\left[|\nabla \psi|^{2}-g^{\prime}\left(\psi_{0}\right) \psi^{2}\right] d x d y=\left(A_{0} \psi, \psi\right)
\end{align*}
$$

where $\psi=(-\Delta)^{-1} \omega$. Recall that $n^{\leq 0}(L)$ and $n \leq 0\left(A_{0}\right)$ denote the maximal dimensions of subspaces where the quadratic forms $\langle L \cdot, \cdot\rangle$ and $\left(A_{0} \cdot, \cdot\right)$ are nonpositive. Let

$$
\left\{\psi_{1}, \cdots, \psi_{l}\right\}, \quad l=n^{\leq 0}\left(A_{0}\right)
$$

be linearly independent eigenfunctions associated to nonpositive eigenvalues of $A_{0}$. Define the space $Y_{1} \subset X$ by

$$
Y_{1}=\left\{\omega \in X \mid \int_{\Omega} \psi_{j}(-\Delta)^{-1} \omega=0,1 \leq j \leq l\right\}
$$

Then for any $\omega \in Y_{1}$, we have

$$
\left(A_{0} \psi, \psi\right) \geq \delta\|\psi\|_{H^{1}}^{2}, \text { for some } \delta>0
$$

So by (11.57), for any $\omega \in Y_{1}$,

$$
\begin{aligned}
\langle L \omega, \omega\rangle & =\varepsilon \iint_{\Omega}\left\{\frac{\omega^{2}}{g^{\prime}\left(\psi_{0}\right)}-|\nabla \psi|^{2}\right\} d x d y+(1-\varepsilon)\langle L \omega, \omega\rangle \\
& \geq \varepsilon \iint_{\Omega}\left\{\frac{\omega^{2}}{g^{\prime}\left(\psi_{0}\right)}-|\nabla \psi|^{2}\right\} d x d y+(1-\varepsilon) \delta\|\psi\|_{H^{1}}^{2} \\
& \geq \varepsilon \iint_{\Omega}\left\{\frac{\omega^{2}}{g^{\prime}\left(\psi_{0}\right)}+|\nabla \psi|^{2}\right\} d x d y
\end{aligned}
$$

by choosing $\varepsilon>0$ such that $(1-\varepsilon) \delta>2 \varepsilon$. Since the positive subspace $Y_{1}$ has co-dimension $n^{\leq 0}\left(A_{0}\right)$, this shows that the assumption (H2) for $L$ on $X$ is satisfied and $n^{\leq 0}(L) \leq n^{\leq 0}\left(A_{0}\right)$.

To prove $n^{\leq 0}(L) \geq n^{\leq 0}\left(A_{0}\right)$, let $\tilde{\omega}_{j}=g^{\prime}\left(\psi_{0}\right) \psi_{j} \in X$ and $\tilde{\psi}_{j}=(-\Delta)^{-1} \tilde{\omega}_{j}$, $j=1, \ldots, l$ and then

$$
\begin{aligned}
\left(A_{0} \psi_{j}, \psi_{j}\right) & =\iint_{\Omega}\left[\left|\nabla \psi_{j}\right|^{2}-g^{\prime}\left(\psi_{0}\right) \psi_{j}^{2}\right] d x d y=\iint_{\Omega}\left[\left|\nabla \psi_{j}\right|^{2}-\frac{\tilde{\omega}_{j}^{2}}{g^{\prime}\left(\psi_{0}\right)}\right] d x d y \\
& =\iint_{\Omega}\left[\left|\nabla \psi_{j}\right|^{2}-2 \tilde{\omega}_{j} \psi_{j}+\frac{\tilde{\omega}_{j}^{2}}{g^{\prime}\left(\psi_{0}\right)}\right] d x d y \\
& =\iint_{\Omega}\left[\left|\nabla \psi_{j}\right|^{2}-2 \nabla \psi_{j} \cdot \nabla \tilde{\psi}_{j}+\frac{\tilde{\omega}_{j}^{2}}{g^{\prime}\left(\psi_{0}\right)}\right] d x d y \\
& \geq \iint_{\Omega}\left[\frac{\tilde{\omega}_{j}^{2}}{g^{\prime}\left(\psi_{0}\right)}-\left|\nabla \tilde{\psi}_{j}\right|^{2}\right] d x d y=\left\langle L \tilde{\omega}_{j}, \tilde{\omega}_{j}\right\rangle
\end{aligned}
$$

and thus $n^{\leq 0}(L) \geq n^{\leq 0}\left(A_{0}\right)$. Combined with above, this implies that $n^{\leq 0}(L)=$ $n^{\leq 0}\left(A_{0}\right)$. Since $\omega \in \operatorname{ker} L$ if and only if $\psi=(-\Delta)^{-1} \omega \in \operatorname{ker} A_{0}$, we obtain $\operatorname{dim} \operatorname{ker} L=\operatorname{dim} \operatorname{ker} A_{0}$ and thus

$$
n^{-}(L)=n^{\leq 0}(L)-\operatorname{dim} \operatorname{ker} L=n^{\leq 0}\left(A_{0}\right)-\operatorname{dim} \operatorname{ker} A_{0}=n^{-}\left(A_{0}\right)
$$

ii) Note that, like $J$, the projection $P$ also commutes with $f\left(\psi_{0}\right)$. for any $f$. Therefore, $\omega \in \overline{R(J)}$ if and only if $P \frac{\omega}{g^{\prime}\left(\psi_{0}\right)}=0$. It implies that

$$
\begin{equation*}
(I-P) L \omega=\frac{\omega}{g^{\prime}\left(\psi_{0}\right)}-(I-P) \psi=\frac{1}{g^{\prime}\left(\psi_{0}\right)} A \psi, \quad \forall \omega \in \overline{R(J)} \tag{11.58}
\end{equation*}
$$

where $\psi=(-\Delta)^{-1} \omega$, and thus

$$
\begin{equation*}
(-\Delta) \operatorname{ker} A=\overline{R(J)} \cap \operatorname{ker}((I-P) L)=\overline{R(J)} \cap \operatorname{ker} J L \tag{11.59}
\end{equation*}
$$

Since

$$
\operatorname{ker}\left(\left.\langle L \cdot, \cdot\rangle\right|_{\overline{R(J)}}\right)=\overline{R(J)} \cap \operatorname{ker} J L
$$

it immediately implies

$$
\begin{equation*}
\operatorname{dim} \operatorname{ker}\left(\left.\langle L \cdot, \cdot\rangle\right|_{\overline{R(J)}}\right)=\operatorname{dim} \operatorname{ker} A \tag{11.60}
\end{equation*}
$$

Suppose $\langle L \cdot, \cdot\rangle$ is degenerate on $\overline{R(J)} /(\overline{R(J)} \cap \operatorname{ker} L)$, namely

$$
\exists \omega_{1} \in \overline{R(J)} \backslash \operatorname{ker} L \text { such that }\left\langle L \omega_{1}, \omega\right\rangle=0, \forall \omega \in \overline{R(J)}
$$

Such $\omega_{1}$ satisfies $0 \neq L \omega_{1} \in \operatorname{ker} J$, or equivalently $(I-P) L \omega_{1}=0$. Therefore, (11.59) implies $A \psi_{1}=0$. Since $A_{0} \psi_{1} \neq 0$ due to $L \omega_{1} \neq 0$, we obtain ker $A \varsubsetneqq \operatorname{ker} A_{0}$. The converse can be proved similarly and the first statement follows.

To prove (11.56), first we notice that for any $\omega \in \overline{R(J)}$,

$$
\begin{aligned}
& \langle L \omega, \omega\rangle=\iint_{\Omega}\left\{\left(\frac{\omega}{\sqrt{g^{\prime}\left(\psi_{0}\right)}}-\psi \sqrt{g^{\prime}\left(\psi_{0}\right)}\right)^{2}-g^{\prime}\left(\psi_{0}\right) \psi^{2}+|\nabla \psi|^{2}\right\} d x d y \\
& =\iint_{\Omega}\left[\left(\frac{\omega}{\sqrt{g^{\prime}\left(\psi_{0}\right)}}-\sqrt{g^{\prime}\left(\psi_{0}\right)}(I-P) \psi\right)^{2}+g^{\prime}\left(\psi_{0}\right)(P \psi)^{2}\right. \\
& \left.\quad-g^{\prime}\left(\psi_{0}\right) \psi^{2}+|\nabla \psi|^{2}\right] d x d y
\end{aligned}
$$

$$
\begin{equation*}
\geq \iint_{\Omega}|\nabla \psi|^{2}-g^{\prime}\left(\psi_{0}\right) \psi^{2}+g^{\prime}\left(\psi_{0}\right)(P \psi)^{2} d x d y=(A \psi, \psi) \tag{11.61}
\end{equation*}
$$

Next, for any $\psi \in H_{0}^{1}$, let

$$
\tilde{\omega}=g^{\prime}\left(\psi_{0}\right)(I-P) \psi \in \overline{R(J)}, \quad \tilde{\psi}=(-\Delta)^{-1} \tilde{\omega}
$$

then

$$
\begin{align*}
(A \psi, \psi) & =\iint_{\Omega}|\nabla \psi|^{2}-g^{\prime}\left(\psi_{0}\right)((I-P) \psi)^{2} d x d y  \tag{11.62}\\
& =\iint_{\Omega}\left[|\nabla \psi|^{2}-\frac{\tilde{\omega}^{2}}{g^{\prime}\left(\psi_{0}\right)}\right] d x d y \\
& =\iint_{\Omega}\left[|\nabla \psi|^{2}-2 \tilde{\omega} \psi+\frac{\tilde{\omega}^{2}}{g^{\prime}\left(\psi_{0}\right)}\right] d x d y \\
& \geq \iint_{\Omega}\left[\frac{\tilde{\omega}^{2}}{g^{\prime}\left(\psi_{0}\right)}-|\nabla \tilde{\psi}|^{2}\right] d x d y=\langle L \tilde{\omega}, \tilde{\omega}\rangle
\end{align*}
$$

From (11.61), (11.62) and (11.60), we get (11.56) as in the proof of i).
By Lemma 11.3 and (iii) of Proposition 2.8, we have
ThEOREM 11.5. Assume $g^{\prime}\left(\psi_{0}\right)>0$ and $\operatorname{ker} A=\{0\}$, then the index formula

$$
\begin{equation*}
k_{r}+2 k_{c}+2 k_{i}^{\leq 0}=n^{-}(A) \tag{11.63}
\end{equation*}
$$

holds. In particular, when $n^{-}(A)$ is odd, there is linear instability; when $A>0$, there is linear stability.

Proof. To apply (iii) of Proposition 2.8 to obtain (11.63), it suffices to verify that a.) $\langle L \cdot, \cdot\rangle$ is non-degenerate on $\operatorname{ker}(J L) / \operatorname{ker} L$, which is satisfied due to Lemma 2.1, Lemma 11.3, and $\{0\}=\operatorname{ker} A \subset \operatorname{ker} A_{0}$; and b.)

$$
\tilde{S} \triangleq \overline{R(J)} \cap(J L)^{-1}(\operatorname{ker} L)=\{0\}
$$

To see the latter, we first note that ker $A=\{0\}$ and (11.59) imply

$$
\overline{R(J)} \cap \operatorname{ker}((I-P) L)=\{0\}
$$

Consequently, if $\omega \in \overline{R(J)} \cap(J L)^{-1}(\operatorname{ker} L)$, then $J L \omega \in R(J) \cap$ ker $L$ must vanish, namely, $L \omega \in \operatorname{ker} J$, and thus $(I-P) L \omega=0$. Again, since $\omega \in \overline{R(J)}$, we obtain $\omega=0$. Therefore, $\tilde{S}=\{0\}$ and (11.63) follows.

The instability of $e^{t J L}$ under the assumption of $n^{-}(A)$ being odd is straightforward from (11.63). Finally suppose $A>0,(11.56)$ and (11.60) imply that $\langle L \cdot, \cdot\rangle$ is uniformly positive definite on $\overline{R(J)} /(\overline{R(J)} \cap \operatorname{ker} L)=\overline{R(J)}$. Therefore, $e^{t J L}$ is stable on the closed invariant subspace $\overline{R(J)}$ and thus its stability follows from the decomposition $X=\operatorname{ker}(J L)+\overline{R(J)}$, which is proved in Proposition 2.8.

By Theorem 2.6, the index formula (11.63) and the fact $i \mathbf{R} \subset \sigma(J L)$ for the linearized Euler equation imply the following.

Corollary 11.5. Under the assumption of Theorem 11.5, when $n^{-}(A)>$ 0 , then there is linear instability or structural instability for $J L$ (in the sense of Theorem 2.6).

In the sense of Theorem 2.6, the structural instability of the linearized Euler equation $\omega_{t}=J L \omega$ means that there exist arbitrarily small bounded perturbations $L_{\#}$ to $L$ such that $J L_{\#}$ has unstable eigenvalues. However, it is not clear that such perturbations can be realized in the context of the Euler equation, such as by considering neighboring steady states along with possible small domain variation.

REMARK 11.7. In [49], it was shown that for general $g \in C^{1}$, when $\operatorname{ker} A=\{0\}$ and $n^{-}(A)$ is odd, there is linear instability. Here, the index formula (11.63) gives more detailed information about the spectrum of the linearized Euler operator.

We give one example satisfying the stability condition $A>0$. Let $\lambda_{0}>0$ be the lowest eigenvalue of $-\Delta$ in $\Omega$ with Dirichlet boundary condition and $\psi_{0}$ be the corresponding eigenfunction. Then $g^{\prime}\left(\psi_{0}\right)=\lambda_{0}$ and it is easy to show that $A>0$.

REMARK 11.8. When the domain $\Omega$ is not simply connected, let $\partial \Omega=\cup_{i=0}^{n} \Gamma_{i}$ consist of outer boundary $\Gamma_{0}$ and $n$ interior boundaries $\Gamma_{1}, \cdots, \Gamma_{n}$. Then the operators $A_{0}, A,-\Delta$ should be defined by using the boundary conditions:

$$
\begin{equation*}
\left.\phi\right|_{\Gamma_{i}} \text { is constant, } \oint_{\Gamma_{i}} \frac{\partial \phi}{\partial n}=0 \text { and } \iint_{\Omega} \phi d x d y=0 . \tag{11.64}
\end{equation*}
$$

The same formula (11.63) is still true. The linearized stream functions satisfying (11.64) represent perturbations preserving the circulations along each $\Gamma_{i}$, which are conserved in the nonlinear evolution.

Below, we consider the case when $\operatorname{ker} A$ is nontrivial. This usually happens when the problem has some symmetry. As an example, we consider the case when $\Omega$ is a channel, that is,

$$
\Omega=\left\{y_{1} \leq y \leq y_{2}, x \text { is } T-\text { periodic }\right\}
$$

The steady stream function $\psi_{0}$ satisfies

$$
\begin{equation*}
-\Delta \psi_{0}=g\left(\psi_{0}\right) \text { in } \Omega \tag{11.65}
\end{equation*}
$$

with boundary conditions $\psi_{0}$ being constants on $\left\{y=y_{i}\right\}, i=1,2$, where $g \in C^{1}$. Define the operators $L, A_{0}, A$ as before with the boundary conditions

$$
\begin{equation*}
\phi \text { is constant on }\left\{y=y_{i}\right\}, \int_{\left\{y=y_{i}\right\}} \frac{\partial \phi}{\partial y} d x=0, i=1,2 \tag{11.66}
\end{equation*}
$$

and $\iint_{\Omega} \phi d x d y=0$. Taking $x$-derivative of equation (11.65), we get

$$
-\Delta \psi_{0, x}=g^{\prime}\left(\psi_{0}\right) \psi_{0, x} \text { in } \Omega
$$

and $\psi_{0, x}$ satisfies the boundary condition (11.66). Thus we have $A_{0} \psi_{0, x}=0$ and

$$
L \omega_{0, x}=L\left(g^{\prime}\left(\psi_{0}\right) \psi_{0, x}\right)=0
$$

Since $\psi_{0, x}=u_{0} \cdot \nabla(-y)$, so $P \psi_{0, x}=0$ and thus $A \psi_{0, x}=A_{0} \psi_{0, x}=0$.
Theorem 11.6. Assume $g^{\prime}\left(\psi_{0}\right)>0$ and $\operatorname{ker} A=\operatorname{span}\left\{\psi_{0, x}\right\}$, then

$$
\exists \omega_{1} \in \overline{R(J)} \text { such that } J L \omega_{1}=\omega_{0, x}
$$

Moreover, if

$$
d=\left\langle L \omega_{1}, \omega_{1}\right\rangle=-\iint_{\Omega} y \omega_{1} d x d y=T\left(\left.\psi_{1}\right|_{y=y_{1}}-\left.\psi_{1}\right|_{y=y_{2}}\right) \neq 0
$$

where $\psi_{1}$ satisfies $-\Delta \psi_{1}=\omega_{1}$ with the boundary condition (11.66), or more explicitly

$$
\psi_{1}=-A^{-1}\left(g^{\prime}\left(\psi_{0}\right)(I-P) y\right)
$$

then we have the index formula

$$
\begin{equation*}
k_{r}+2 k_{c}+2 k_{i}^{\leq 0}=n^{-}(A)-n^{-}(d) \tag{11.67}
\end{equation*}
$$

where $n^{-}(d)=1$ if $d<0$ and $n^{-}(d)=0$ if $d>0$.
Proof. Our assumption implies ker $A \subset \operatorname{ker} A_{0}$ and thus $\langle L \cdot, \cdot\rangle$ is non-degenerate on $\overline{R(J)} /(\overline{R(J)} \cap$ ker $L)$ by Proposition 2.8. To apply index formula (2.20), we need to obtain the non-degeneracy of $\langle L \cdot, \cdot\rangle$ on $\tilde{S} /(\tilde{S} \cap \operatorname{ker} L)$ and compute $n^{\leq 0}\left(\left.L\right|_{\tilde{S} /(\tilde{S} \cap \operatorname{ker} L)}\right)$, where

$$
\tilde{S}=\overline{R(J)} \cap(J L)^{-1}(\operatorname{ker} L)=\overline{R(J)} \cap(J L)^{-1}(\operatorname{ker} L \cap \overline{R(J)})
$$

Since $P \psi_{0, x}=0$ implies $\omega_{0, x}=\frac{\psi_{0, x}}{g^{\prime}\left(\psi_{0}\right)} \in \overline{R(J)}$. From (11.59) and our assumption on $\operatorname{ker} A$, we have

$$
\operatorname{span}\left\{\omega_{0, x}\right\} \subset \overline{R(J)} \cap \operatorname{ker} L \subset \overline{R(J)} \cap \operatorname{ker}((I-P) L)=\Delta \operatorname{ker} A=\operatorname{span}\left\{\omega_{0, x}\right\}
$$

which yields

$$
\overline{R(J)} \cap \operatorname{ker} L=\operatorname{span}\left\{\omega_{0, x}\right\}
$$

By the definition of $\tilde{S}, \omega \in \tilde{S}$ if and only if there exist $\omega \in \overline{R(J)}$ and $a \in \mathbf{R}$ such that

$$
a \omega_{0, x}=J L \omega=J(I-P) L \omega
$$

Since $\omega_{0, x}=-J y=-J(I-P) y$, we obtain equivalently $(I-P)(L \omega+a y)=0$.
From (11.58), it follows that

$$
\omega \in \tilde{S} \Longleftrightarrow \omega \in \overline{R(J)} \text { and } A \psi=-a g^{\prime}\left(\psi_{0}\right)(I-P) y, \text { where }-\Delta \psi=\omega
$$

Note $\operatorname{ker} A=\operatorname{span}\left\{\psi_{0, x}\right\}$ and

$$
\left\langle g^{\prime}\left(\psi_{0}\right)(I-P) y, \psi_{0, x}\right\rangle=\iint_{\Omega} y g^{\prime}\left(\psi_{0}\right) \psi_{0, x} d x d y=0
$$

There exists a stream function $\psi_{1}$ satisfying

$$
A \psi_{1}=-g^{\prime}\left(\psi_{0}\right)(I-P) y
$$

which implies $\omega_{1}=-\Delta \psi_{1} \in \overline{R(J)}$ and $J L \omega_{1}=\omega_{0, x}$. Namely $\omega_{1} \in \tilde{S}$ and $\tilde{S}=$ $\operatorname{span}\left\{\omega_{0, x}, \omega_{1}\right\}$. One may compute

$$
\begin{aligned}
d & =\left\langle L \omega_{1}, \omega_{1}\right\rangle=\left\langle(I-P) L \omega_{1}, \omega_{1}\right\rangle=\left\langle-(I-P) y, \omega_{1}\right\rangle \\
& =-\left\langle y, \omega_{1}\right\rangle=T\left(\left.\psi_{1}\right|_{y=y_{1}}-\left.\psi_{1}\right|_{y=y_{2}}\right)
\end{aligned}
$$

where the last equal sign follows from integration by parts. If $d \neq 0$, then the desired index formula follows from (iii) of Proposition 2.8.

Similar to Corollary 11.5 (and the comments immediately thereafter), we have
Corollary 11.6. Under the assumption of Theorem 11.6, when $n^{-}(A)-$ $n^{-}(d)>0$, then there is linear instability or structural instability for $J L$.

As another application of the Hamiltonian structure of the linearized Euler equation, we consider the inviscid damping of a stable steady flow. Assume $g^{\prime}\left(\psi_{0}\right)>0$ and $A>0$, then by Theorem 11.5 , the steady flow is linearly stable in the $L^{2}$ norm of vorticity. There is no time decay in $\|\omega\|_{L^{2}}$. However, the linear decay in the velocity norm $\|u\|_{L^{2}}$ is possible due to the mixing of the vorticity. For example, see [56] for the linear damping near Couette flow $(y, 0)$ in a channel. Here, we give a weak form of the linear decay for general stable steady flows.

Theorem 11.7. Assume $g^{\prime}\left(\psi_{0}\right)>0$ and $A>0$. For $\omega(0) \in \overline{R(J)}$, let $\omega(t) \in$ $\overline{R(J)}$ be the solution of the linearized Euler equation (11.54). Then
(i) When $T \rightarrow \infty, \frac{1}{T} \int_{0}^{T} \omega(t) d t \rightarrow 0$ strongly in $L^{2}$.
(ii) If there is no embedded imaginary eigenvalue of $J L$ on $\overline{R(J)}$, then for any compact operator $C$ in $L^{2}$, we have

$$
\begin{equation*}
\frac{1}{T} \int_{0}^{T}\|C \omega(t)\|_{L^{2}}^{2} d t \rightarrow 0, \text { when } T \rightarrow \infty \tag{11.68}
\end{equation*}
$$

In particular, for the velocity $u=\operatorname{curl}^{-1} \omega$,

$$
\begin{equation*}
\frac{1}{T} \int_{0}^{T}\|u(t)\|_{L^{2}}^{2} d t \rightarrow 0, \text { when } T \rightarrow \infty \tag{11.69}
\end{equation*}
$$

Proof. By Lemma 11.3, $\left.L\right|_{\overline{R(J)}}>0$. Since $\overline{R(J)}$ is an invariant subspace of $J L$, we can consider the operator $J L$ in $\overline{R(J)}$. Define the inner product $[\cdot, \cdot]=\langle L \cdot, \cdot\rangle$ on $\overline{R(J)}$, then the norm in $[\cdot, \cdot]$ is equivalent to the $L^{2}$ norm. As noted before, the operator $\left.J L\right|_{\overline{R(J)}}$ is anti-self-adjoint with respect to the inner product $[\cdot, \cdot]$.
(i) By the mean ergodic convergence of unitary operators ([73])

$$
\lim _{T \rightarrow \infty} \frac{1}{T} \int_{0}^{T} \omega(t) d t=\lim _{T \rightarrow \infty} \frac{1}{T} \int_{0}^{T} e^{\left.t J L\right|_{\overline{R(J)}}} \omega(0) d t=P_{0} \omega(0)
$$

in $L^{2}$, where $P_{0}$ is the projection operator from $\overline{R(J)}$ to ker $\left.J L\right|_{\overline{R(J)}}$ orthogonal with respect to $[\cdot, \cdot]$. Since $\operatorname{ker} A=\{0\}$, by Lemma 11.3 and (11.58) in particular, $\left.\operatorname{ker} J L\right|_{\overline{R(J)}}=\{0\}$ and thus $P_{0} \omega(0)=0$.
(ii) If $J L$ has no embedded imaginary eigenvalue, then (11.68) follows directly by the RAGE theorem ([22]), again by using the anti-self-adjoint property of $\left.J L\right|_{\overline{R(J)}}$. The conclusion (11.69) follows by choosing the compact operator $C=\operatorname{curl}^{-1}$.

Remark 11.9. Assuming $A>0$, from the proof of Theorem 11.5, the subspace $\tilde{S}$ defined in Proposition 2.8 is trivial. By Proposition 2.8, there is a direct sum decomposition $L^{2}=\operatorname{ker}(J L) \oplus \overline{R(J)}$ invariant under $J L$. In fact $\operatorname{ker}(J L)$ corresponds to the steady solution of the linearized Euler equation. So above Lemma shows that for any initial data in $L^{2}$, in the time averaged limit, the solution of the linearized Euler equation converges to a steady solution. This is a weak form of inviscid damping.

A stable example satisfying the assumption $A>0$ in Theorem 11.7 is given in Remark 11.7. Below, we consider two examples of stable shear flows. First, we consider the Poisseulle flow $U(y)=y^{2}$ in a $2 \pi$-periodic channel $\{-1<y<1\}$. The linearized Euler equation becomes

$$
\partial_{t} \omega+y^{2} \partial_{x} \omega+2 \partial_{x} \psi=0
$$

Consider the subspace of non-shear vorticities with a weighted $L^{2}$ norm

$$
X_{1}=\left\{\omega=\sum_{k \in \mathbf{Z}, k \neq 0} e^{i k x} \omega_{k}(y),\|\omega\|_{X_{1}}^{2}=\sum_{k \in \mathbf{Z}, k \neq 0}\left\|y \omega_{k}\right\|_{L^{2}}^{2}<\infty\right\}
$$

Define $J=-\partial_{x}$ and $L=y^{2}+2(-\Delta)^{-1}$. Then $L$ is uniformly positive on $X_{1}$.
Second, consider the Kolmogorov flow $U(y)=\sin y$ in a torus $T^{2}=S_{\frac{2 \pi}{\alpha}} \times S_{2 \pi}$ with $\alpha>1$. Here $\alpha>1$ is the sharp stability condition since the shear flow is unstable when $\alpha<1$. The linearized equation is

$$
\partial_{t} \omega+\sin y \partial_{x}(\omega-\psi)=0
$$

Let $J=\sin y \partial_{x}$ and $L=1-(-\Delta)^{-1}$. Then $L$ is uniformly positive on

$$
X_{2}=\left\{\omega=\sum_{k \in \mathbf{Z}, k \neq 0} e^{i k x} \omega_{k}(y), \omega \in L^{2}\right\}
$$

when $\alpha>1$. It can be shown ([54]) that for above two examples, the linearized Euler operator has no embedded eigenvalues. Therefore, Theorem 11.7 (ii) is true for the above two shear flows in $X_{1}$ and $X_{2}$ respectively. In particular, if we choose $C$ to be $P_{N}$, the projection operator to the first $N$ Fourier modes (in $x$ ), then

$$
\begin{equation*}
\frac{1}{T} \int_{0}^{T}\left\|P_{N} \omega(t)\right\|_{L^{2}}^{2} d t \rightarrow 0, \text { when } T \rightarrow \infty \tag{11.70}
\end{equation*}
$$

This shows that in the time averaged sense, the low frequency parts of $\omega$ tends to zero. This observation was used to prove ([54]) the metastability of Kolmogorov flows. In the fluid literature (see e.g. [71]), for 2D turbulence a dual cascade was known that energy moves to low frequency end and the enstrophy ( $\int \omega^{2} d x$ ) moves to the high frequency end. The result (11.70) can be seen as a justification of such physical intuition in a weak sense.

REmARK 11.10. Two classes of shear flows generalizing the above two examples are studied in [54]. The linear inviscid damping in the sense of (11.69) is proved for stable shear flows and on the center space for the unstable shear flows, when $\omega(0) \in L^{2}$ is non-shear. Recently, for monotone and certain symmetric shear flows, more explicit linear decay estimates of the velocity were obtained in [77, 74, 75] for more regular initial data (e.g. $\omega(0) \in H^{1}$ or $H^{2}$ ).

In [55], the stability of shear flows under Coriolis forces is studied. By using the instability index Theorem 2.3, the sharp stability condition for a class of shear flows can be obtained. Then the linear damping as in the above sense is proved for non-shear $\omega(0) \in L^{2}$.

### 11.6. Stability of traveling waves of 2D NLS

In this section, we consider the nonlinear Schrödinger equation (NLS)

$$
\begin{equation*}
i \frac{\partial u}{\partial t}+\Delta u+F\left(|u|^{2}\right) u=0, \quad u=u_{1}+i u_{2}: \mathbf{R} \times \mathbf{R}^{2} \rightarrow \mathbf{C} \tag{11.71}
\end{equation*}
$$

In particular, we assume that the nonlinearity $F(s)$ satisfies

$$
\begin{equation*}
F \in C^{2}, \quad F(1)=0, \quad F^{\prime}(1)<0 \tag{11.72}
\end{equation*}
$$

Important well-known equations of this type are Gross-Pitaevskii (GP) equation with $F(s)=1-s$ and the cubic-quintic NLS with $F(s)=-\alpha_{1}+\alpha_{3} s-\alpha_{5} s^{2}$, where $\alpha_{1}, \alpha_{3}$ and $\alpha_{5}$ are positive constants. Assume $s=1$ is a local minimal point of $F$, it is natural to consider solutions $u(t, x)$ satisfying the following boundary condition in some appropriate sense

$$
\begin{equation*}
|u| \rightarrow 1 \text { as }|x| \rightarrow \infty \tag{11.73}
\end{equation*}
$$

After normalization, we can assume that $u \rightarrow 1$ when $|x| \rightarrow \infty$ in some weak sense such as $u-1$ being approximable by Schwartz class functions in certain Sobolev norms. The equation (11.71) has the conserved energy and momentum functionals

$$
\begin{gathered}
E(u)=\frac{1}{2} \int_{\mathbf{R}^{2}}\left(|\nabla u|^{2}+V\left(|u|^{2}\right)\right) d x \\
\vec{P}(u)=\left(P_{1}(u), P_{2}(u)\right)=\frac{1}{2} \int_{\mathbf{R}^{2}}\langle\nabla u, i(u-1)\rangle d x=\int_{\mathbf{R}^{2}}\left(u_{1}-1\right) \nabla u_{2} d x
\end{gathered}
$$

where $V(s)=\int_{s}^{1} F(\tau) d \tau$. We also denote the first component of $\vec{P}(u)$ by

$$
P(u)=\frac{1}{2} \int_{\mathbf{R}^{2}}\left\langle\partial_{x_{1}} u, i(u-1)\right\rangle d x=\int_{\mathbf{R}^{2}}\left(u_{1}-1\right) \partial_{x_{1}} u_{2} d x
$$

A traveling wave (without loss of generality, in $x_{1}$-direction) of (11.71) with wave speed $c \in(0, \sqrt{2})$ is a solution in the form of $u=U_{c}\left(x_{1}-c t, x_{2}\right)$, where $U_{c}$ satisfies the elliptic equation

$$
\begin{equation*}
-i c \partial_{x_{1}} U_{c}+\Delta U_{c}+F\left(\left|U_{c}\right|^{2}\right) U_{c}=0 \tag{11.74}
\end{equation*}
$$

with the boundary condition $U_{c} \rightarrow 1$ when $|x| \rightarrow \infty$ in the sense $U_{c}-1 \in \dot{H}^{1}$. Here, $\sqrt{2}$ is the sound speed and when $c \geq \sqrt{2}$, in general the traveling waves do not exist (see e.g. [61]). Formally, $U_{c}$ is a critical point of $E-c P$. Our goal is to understand the linear stability/instability of such a traveling wave, namely, the
evolution of the linearized equation of (11.71) at $U_{c}=u_{c}+i v_{c}$ put in the moving frame $x_{1} \rightarrow x_{1}-c t, x_{2} \rightarrow x_{2}$ :

$$
\begin{equation*}
u_{t}=J L_{c} u, \quad u=\left(u_{1}, u_{2}\right)^{T} \rightarrow 0 \text { as }|x| \rightarrow \infty \tag{11.75}
\end{equation*}
$$

where $J=\left(\begin{array}{cc}0 & 1 \\ -1 & 0\end{array}\right)$ and
$L_{c}:=\left(\begin{array}{cc}-\Delta-F\left(\left|U_{c}\right|^{2}\right)-2 F^{\prime}\left(\left|U_{c}\right|^{2}\right) u_{c}^{2} & -c \partial_{x_{1}}-2 F^{\prime}\left(\left|U_{c}\right|^{2}\right) u_{c} v_{c} \\ c \partial_{x_{1}}-2 F^{\prime}\left(\left|U_{c}\right|^{2}\right) u_{c} v_{c} & -\Delta-F\left(\left|U_{c}\right|^{2}\right)-2 F^{\prime}\left(\left|U_{c}\right|^{2}\right) v_{c}^{2}\end{array}\right)$.
Through $L^{2}$ duality, $L_{c}$ generates the quadratic form

$$
\begin{align*}
\left\langle L_{c} u, v\right\rangle=\int_{\mathbf{R}^{2}}\left\{\nabla u \cdot \nabla v+c\left(v_{1 x_{1}} u_{2}\right.\right. & \left.-u_{1} v_{2 x_{1}}\right)-F\left(\left|U_{c}\right|^{2}\right) u \cdot v \\
& \left.-2 F^{\prime}\left(\left|U_{c}\right|^{2}\right)\left(U_{c} \cdot u\right)\left(U_{c} \cdot v\right)\right\} d x \tag{11.76}
\end{align*}
$$

where $u \cdot v=\operatorname{Re}(u \bar{v})$.
For the purpose of studying the linearized equation (11.75), we make the following assumptions:
(NLS-1) $U_{c}-1 \in H^{1} \times \dot{H}^{1}$ satisfies (11.73) and $\left|U_{c}\right|_{C^{1}\left(\mathbf{R}^{2}\right)}<\infty$.
(NLS-2) Let $\Gamma$ be the collection of subspaces $S \subset H^{1}\left(\mathbf{R}^{2}\right) \times H^{1}\left(\mathbf{R}^{2}\right)$ such that $\left\langle L_{c} u, u\right\rangle<0$ for all $0 \neq u \in S$, then

$$
\max \{\operatorname{dim} S \mid S \in \Gamma\}=n^{-}\left(L_{c}\right)<\infty
$$

The above (NLS-1) is a natural regularity assumption. For any given traveling wave of (11.71), it is probably not so straightforward to verify (NLS-2). This, however, would be a direct consequence if $U_{c}$ is obtained through a constrained variational approach related to energy and momentum, which is often the case. For example, in [19] [20], the 2D traveling waves of (11.74) were constructed by minimizing the functional $E(u)-c P(u)$ subject to a constraint $P(u)=p$ or $E_{k i n}(u)=\int|\nabla u|^{2} d x=k$, for general nonlinearity $F$. The variational problem of minimizing $E(u)-c P(u)$ subject to fixed $P(u)$ was also studied in [7] to construct 2D traveling waves of GP equation. Since these 2D traveling waves $U_{c}$ were minimizers of $E(u)-c P(u)$ subject to one constraint, it can be shown that $n^{-}\left(L_{c}\right) \leq 1$ (see e.g. the proof of Lemma 2.7 of [53]). Here, we note that $U_{c}$ is a critical point of $E(u)-c P(u)$ and $L_{c}=E^{\prime \prime}\left(U_{c}\right)-c P^{\prime \prime}\left(U_{c}\right)$.

To study the quadratic form $\left\langle L_{c} \cdot, \cdot\right\rangle$, obviously one may take $X=H^{1}\left(\mathbf{R}^{2}\right) \times$ $H^{1}\left(\mathbf{R}^{2}\right)$. On the one hand, the above assumptions ensure that $L_{c}: X \rightarrow X^{*}=$ $H^{-1} \times H^{-1}$ is bounded, satisfies $L_{c}^{*}=L_{c}$, and has $n^{-}\left(L_{c}\right)$ negative dimensions. On the other hand, it is easy to see that $J: X^{*} \rightarrow X$ is unbounded, but has a dense domain $H^{1} \times H^{1} \subset X^{*}=H^{-1} \times H^{-1}$, and satisfies $J^{*}=-J$. However, as the boundary condition (11.73) does not provide enough control of $|u|^{2}$ near $|x|=\infty$ in $\left\langle L_{c} u, u\right\rangle$, it is not clear that (H2.b) can be satisfied by any decomposition.

For (11.71) considered on $\mathbf{R}^{N}, N \geq 3$, as in [53], it would be possible to work on $X=H^{1} \times \dot{H}^{1}$, where $u_{1} \in H^{1}$ and $u_{2} \in \dot{H}^{1}$, and verify assumptions (H1-3) for $J$ and $L_{c}$ based on the following two observations. Firstly, in such higher dimensions, the Gagliardo-Nirenberg inequality implies that $\dot{H}^{1}$ functions decay at $x=\infty$ in the $L^{p}$ sense. Therefore, we may reasonably strengthen the boundary condition (11.73) to $U_{c} \rightarrow 1$ as $|x| \rightarrow \infty$. Consequently the 'principle part' in $\left\langle L_{c} u, u\right\rangle$ provides the control on the $H^{1} \times \dot{H}^{1}$ norm of $u$. Secondly, there are indications that $U_{c}$ decays
like $u_{c}-1=O\left(|x|^{-N}\right)$ and $v_{c}=O\left(|x|^{1-N}\right)$ as in the case proved for the (GP) equation in [6]. Along with the Hardy inequality, this allows us to control those terms in (11.76) with vanishing variable coefficients by the $H^{1} \times \dot{H}^{1}$ norm of $u$. See [53] for more details.

The situation is much worse on $\mathbf{R}^{2}$ unfortunately since both of the above key observations break down on $\mathbf{R}^{2}$. To overcome these difficulties, our idea is to study the stability of the linearized equation (11.75) on some space roughly between $H^{1} \times H^{1}$ and $\dot{H}^{1} \times \dot{H}^{1}$ defined according to the properties of $L_{c}$ by applying Theorem 2.7.

Let $X=H^{1} \times H^{1}$ for (11.75) and define $Q_{0}, Q_{1}: X \rightarrow X^{*}$ as

$$
\left\langle Q_{0} u, v\right\rangle=\operatorname{Re} \int_{\mathbf{R}^{2}} u \bar{v} d x, \quad\left\langle Q_{1} u, v\right\rangle=\operatorname{Re} \int_{\mathbf{R}^{2}}\left(u_{x_{1}} \bar{v}_{x_{1}}+u_{x_{2}} \bar{v}_{x_{2}}\right) d x
$$

namely, the $L^{2}$ and $\dot{H}^{1}$ duality, respectively, which satisfy (B1) in Section 2.6. Let $\mathbb{J}: X \rightarrow X$ be $\mathbb{J}=\left(\begin{array}{cc}0 & 1 \\ -1 & 0\end{array}\right)$. Clearly, $\mathbb{J}$ satisfies (B2) and the unbounded operator $J=\mathbb{J} Q_{0}^{-1}: X^{*} \rightarrow X$ has the same matrix representation through the $L^{2}$ duality. As $L_{c}-Q_{1}$ consists of terms of at most one order of derivative, it satisfies (B3). From (NLS-2), there exists a subspace $S \subset X$ such that $\operatorname{dim} S=n^{-}\left(L_{c}\right)$ and $L_{c}$ is negative definite on $S$. By a slight perturbation, e.g. applying the mollifier to a basis of $S$, we obtain a subspace $X_{-} \subset H^{3} \times H^{3}$ such that $\operatorname{dim} X_{-}=n^{-}\left(L_{c}\right)$ and $L_{c}$ is negative definite on $X_{-}$. Let

$$
X_{\geq 0}=X_{-}^{\perp_{L_{c}}}=\left\{u \in X \mid\left\langle L_{c} v, u\right\rangle=0\right\} \supset \operatorname{ker} L_{c}
$$

and

$$
X_{+}=\left\{u \in X_{\geq 0} \mid \int_{\mathbf{R}^{2}} u \cdot v d x=0, \forall v \in \operatorname{ker} L\right\}
$$

Since $\operatorname{dim} X_{-}<\infty$ and $L_{c}$ is negative definite on $X_{-}$, from Lemma 12.2 where (H2.b) is not necessary (see Remark 12.1), we have $X=X_{-} \oplus X_{\geq 0}$. It is obvious $X_{\geq 0}=X_{+} \oplus \operatorname{ker} L_{c}$ and the decomposition $X=X_{-} \oplus \operatorname{ker} L_{c} \oplus X_{+}$is $L_{c}$-orthogonal. From (NLS-2) and the definition of $X_{ \pm},\left\langle L_{c} u, u\right\rangle$ is (not necessarily uniformly) positive on $X_{+}$and thus ( $\mathbf{B 4}$ ) is satisfied. Finally, one may compute from the construction that

$$
\operatorname{ker} i_{X_{+}}^{*}=Q_{0}\left(\operatorname{ker} L_{c}\right) \oplus L_{c}\left(X_{-}\right)
$$

Since we take $X_{-} \subset H^{3} \times H^{3}$ and $\left|U_{c}\right|_{C^{1}}<\infty,(\mathbf{B 5})$ is also satisfied. From Theorem 2.7, there exists a function space $Y$ roughly between $X=H^{1} \times H^{1}$ and $\dot{H}^{1} \times \dot{H}^{1}$, an extension $L_{c, Y}: Y \rightarrow Y^{*}$ of $L_{c}$, and the restriction

$$
J_{Y}: Y^{*} \supset D\left(J_{Y}\right) \rightarrow Y
$$

of $J$, such that $\left(Y, L_{c, Y}, J_{Y}\right)$ satisfies assumption (H1-3). Therefore, all our main results apply to the linearized NLS (11.75) on $Y$.

In the rest of this section, we assume, for some $c_{0}>0$,
(NLS) There exists a $C^{1}$ curve of traveling waves for $c$ near $c_{0}$ satisfying (NLS-1) such that $n^{-}\left(L_{c_{0}}\right) \leq 1$ and (NLS-2) is satisfied for $c=c_{0}$.
As mentioned in the above, $n^{-}\left(L_{c_{0}}\right) \leq 1$ is satisfied if $U_{c_{0}}$ is constructed as minimizers of $E-c_{0} P$ subject to one constraint such as fixed $P(u)$ or $E_{k i n}(u)$. We shall apply Theorem 2.3 to study the linearized equation (11.75) on $Y$. In order
to estimate $k_{0}^{\leq 0}$ in the counting formula (2.13), differentiating (11.74) in $x_{i}$ and we get ker $L_{c_{0}} \supset\left\{\partial_{x_{i}} U_{c_{0}}, i=1,2\right\}$. Moreover, differentiating (11.74) in $c$, we have

$$
\left.L_{c_{0}} \partial_{c} U_{c}\right|_{c_{0}}=P^{\prime}\left(U_{c_{0}}\right)=J^{-1} \partial_{x_{1}} U_{c_{0}}
$$

and thus $\left.J L_{c_{0}} \partial_{c} U_{c}\right|_{c_{0}} \in \operatorname{ker} L_{c_{0}}$. Since

$$
\left\langle\left. L_{c_{0}} \partial_{c} U_{c}\right|_{c_{0}},\left.\partial_{c} U_{c}\right|_{c_{0}}\right\rangle=\left.\frac{d P\left(U_{c}\right)}{d c}\right|_{c_{0}}
$$

by Proposition 2.7, we have $k_{0}^{\leq 0} \geq 1$ when $\left.\frac{d P\left(U_{c}\right)}{d c}\right|_{c_{0}} \leq 0$ and in this case $U_{c}$ is spectrally stable by (2.13).

The traveling waves constructed in the literature ([7] [19] [20]) are even in $x_{2}$, that is, of the form $U_{c}\left(x_{1},\left|x_{2}\right|\right)$. Thus, we can consider odd and even perturbations (in $x_{2}$ ) respectively. We consider the even perturbations, that is, in the space $Y_{e}=\left\{u \in Y \mid u\right.$ is even in $\left.x_{2}\right\}$. For traveling waves as constrained minimizers of $E-c P$, in general it can be shown that there is at least one even negative direction of $\left\langle L_{c} \cdot, \cdot\right\rangle$, which then implies $n^{-}\left(\left.L_{c}\right|_{Y_{e}}\right)=1$. Such a symmetry preserving negative direction of $L_{c}$ was constructed in [53] for the 3D case. For the 2D case, an even negative direction could be constructed by refining the Derrick type arguments used in [39]. More specifically, one can consider a scaled traveling wave $U^{a, b}=$ $U_{c_{0}}\left(a x_{1}, b x_{2}\right)$ and choose a family of parameters $a(s), b(s)$ near 1 with $a(0)=$ $b(0)=1$ such that

$$
\left(E-c_{0} P\right)\left(U^{a, b}\right)<\left(E-c_{0} P\right)\left(U_{c_{0}}\right),
$$

from which an even negative direction $\left.\frac{d}{d s} U^{a(s), b(s)}\right|_{s=0}$ may be obtained. If in addition to the condition $n^{-}\left(\left.L_{c}\right|_{Y_{e}}\right)=1$, we assume that $\partial_{x_{1}} U_{c}$ is the only even kernel of $L_{c}$, then by Theorem 2.3 and Proposition 2.7, there is linear instability in case $\left.\frac{d P\left(U_{c}\right)}{d c}\right|_{c_{0}}>0$. We summarize above discussions in the following theorem.

Theorem 11.8. (i) Assuming (NLS), the 2D traveling wave $U_{c_{0}}$ is spectrally stable if $\left.\frac{d P\left(U_{c}\right)}{d c}\right|_{c_{0}} \leq 0$.
(ii) If we further assume that $U_{c_{0}}$ is even in $x_{2}$ and there exists $v \in Y_{e}$ in the negative direction of $L_{c_{0}}$ and ker $L_{c_{0}} \cap Y_{e}=\operatorname{span}\left\{\partial_{x_{1}} U_{c_{0}}\right\}$, then $\left.\frac{d P\left(U_{c}\right)}{d c}\right|_{c_{0}}>0$ implies linear instability of $U_{c_{0}}$.

For the GP equation, by numerical computations ([39]) $d P / d c<0$ is true for the whole solitary wave branch. Thus 2D traveling waves of GP are expected to be linearly stable. In [19], the orbital stability of these GP traveling waves was obtained by showing concentration compactness of the constrained minimizing sequence, under the assumption of local uniqueness of minimizers. The transversal instability of 2D traveling waves of GP to 3D perturbation was proved in [53]. For general nonlinear term $F$ such as cubic-quintic type, it is possible that there is an unstable branch of 2D traveling waves with $d P / d c>0$. See the numerical examples given in [20].

Lastly, as a corollary of Theorems 2.3 and 11.8, we prove that the traveling waves $U_{c_{0}}$ have positive momentum $P\left(U_{c_{0}}\right)$.

Corollary 11.7. Under the assumptions in both (i) and (ii) of Theorem 11.8, except for the signs of $\left.\frac{d P\left(U_{c}\right)}{d c}\right|_{c_{0}}$, we have $P\left(U_{c_{0}}\right)>0$.

Proof. First, we find $v_{2}$ such that $L_{c_{0}} v_{2}=J^{-1} \partial_{x_{2}} U_{c_{0}}$. Consider traveling waves $U_{\vec{c}}(\vec{x}-\vec{c} t)$ with velocity vector $\vec{c}=\left(c_{1}, c_{2}\right)$ and $|\vec{c}|=c \in(0, \sqrt{2})$, which satisfies

$$
\begin{equation*}
-J \vec{c} \cdot \nabla U_{\vec{c}}+\Delta U_{\vec{c}}+F\left(\left|U_{\vec{c}}\right|^{2}\right) U_{\vec{c}}=0 \tag{11.77}
\end{equation*}
$$

Let

$$
Q=\frac{1}{|\vec{c}|}\left(\begin{array}{cc}
c_{1} & c_{2} \\
-c_{2} & c_{1}
\end{array}\right)
$$

be the rotating matrix which transforms $\vec{c}$ to $(c, 0)$, then it is easy to check that $U_{\vec{c}}(\vec{x})=U_{c}(Q \vec{x})$ is a solution of (11.77) and

$$
\vec{P}\left(U_{\vec{c}}\right)=Q^{T} \vec{P}\left(U_{c}\right)=P\left(U_{c}\right) \frac{\vec{c}}{c}
$$

where we use $\vec{P}\left(U_{c}\right)=P\left(U_{c}\right)(1,0)^{T}$ which is due to the evenness of $U_{c}$ in $x_{2}$. Differentiating (11.77) in $c_{2}$ and then evaluating at $\left(c_{0}, 0\right)$, we get

$$
\left.L_{c_{0}} \partial_{c_{2}} U_{\vec{c}}\right|_{\left(c_{0}, 0\right)}=J^{-1} \partial_{x_{2}} U_{c_{0}} .
$$

Thus we can choose $v_{2}=\left.\partial_{c_{2}} U_{\vec{c}}\right|_{\left(c_{0}, 0\right)}$ and

$$
\left\langle L_{c_{0}} v_{2}, v_{2}\right\rangle=\left.\partial_{c_{2}} P_{2}\left(U_{\vec{c}}\right)\right|_{\left(c_{0}, 0\right)}=\left.\partial_{c_{2}}\left(P\left(U_{c}\right) \frac{c_{2}}{c}\right)\right|_{\left(c_{0}, 0\right)}=\frac{P\left(U_{c_{0}}\right)}{c_{0}}
$$

Denote $v_{1}=\left.\partial_{c} U_{c}\right|_{c_{0}}$ and recall that

$$
L_{c_{0}} v_{1}=J^{-1} \partial_{x_{1}} U_{c_{0}},\left\langle L_{c_{0}} v_{1}, v_{1}\right\rangle=\left.\frac{d P\left(U_{c}\right)}{d c}\right|_{c_{0}}
$$

Also, by using the evenness of $U_{c_{0}}$ in $x_{2}$, we get

$$
\left\langle L_{c_{0}} v_{2}, v_{1}\right\rangle=\left\langle J^{-1} \partial_{x_{2}} U_{c_{0}},\left.\partial_{c} U_{c}\right|_{c_{0}}\right\rangle=0
$$

and thus

$$
\left.\left\langle L_{c_{0}} \cdot, \cdot\right\rangle\right|_{\operatorname{span}\left\{v_{1}, v_{2}\right\}}=\left(\begin{array}{cc}
\left.\frac{d P\left(U_{c}\right)}{d c}\right|_{c_{0}} & 0 \\
0 & \frac{P\left(U_{c_{0}}\right)}{c_{0}}
\end{array}\right) .
$$

Since

$$
n^{\leq 0}\left(\left.L_{c_{0}}\right|_{\text {span }\left\{v_{1}, v_{2}\right\}}\right) \leq k_{0}^{\leq 0}\left(L_{c_{0}}\right) \leq n^{-}\left(L_{c_{0}}\right) \leq 1
$$

when $\left.\frac{d P\left(U_{c}\right)}{d c}\right|_{c_{0}} \leq 0$, we must have $P\left(U_{c_{0}}\right)>0$. When $\left.\frac{d P\left(U_{c}\right)}{d c}\right|_{c_{0}}>0$ and with the assumptions of Theorem 11.8 (ii), $U_{c_{0}}$ is linearly unstable, which again implies that $P\left(U_{c_{0}}\right)>0$. Since otherwise $P\left(U_{c_{0}}\right) \leq 0$, then $k_{0}^{\leq 0}\left(L_{c_{0}}\right) \geq 1$ and by Theorem 2.3, $U_{c_{0}}$ is linearly stable, a contradiction.

REmark 11.11. For 2D traveling wave solution $U_{c}$ satisfying (11.74), one can prove the identity

$$
\begin{equation*}
c P\left(U_{c}\right)=2 \int_{\mathbf{R}^{2}} V\left(\left|U_{c}\right|\right)^{2} d x \tag{11.78}
\end{equation*}
$$

by using energy conservation and virial identity (see [20] for general $F$ and [39] for $G P)$. So for $F$ such that $V$ is nonnegative (such as $G P$ ), we have $P\left(U_{c}\right)>0$ from (11.78). However, when $V$ also takes negative values (such as cubic-quintic), then one can not conclude the sign of $P\left(U_{c}\right)$ from (11.78). By using the index counting, above Corollary 11.7 shows that $P\left(U_{c}\right)>0$ is true for any nonlinear term $F$ under the assumptions there.

Consider axial symmetric 3D traveling waves $U_{c}=\left(x_{1},\left|x^{\perp}\right|\right)$ which are constrained energy-momentum minimizers, as constructed in [60]. We can also prove that $P\left(U_{c}\right)>0$ by the same arguments as in Corollary 11.7. Actually, the argument for 3D is much simpler than 2D and does not need the additional assumptions on $\operatorname{ker} L_{c}$. Let

$$
v_{1}=\partial_{c} U_{c}, v_{j}=\left.\partial_{c_{j}} U_{\vec{c}}\right|_{(c, 0,0)}, j=2,3
$$

where $\vec{c}=\left(c_{1}, c_{2}, c_{3}\right)$ with $|\vec{c}|=c \in(0, \sqrt{2})$ and $U_{\vec{c}}$ is the traveling wave with the velocity vector $\vec{c}$. Then we can compute in a similar way that

$$
\left.\left\langle L_{c} \cdot, \cdot\right\rangle\right|_{\operatorname{span}\left\{v_{1}, v_{2}, v_{3}\right\}}=\left(\begin{array}{ccc}
\frac{d P\left(U_{c}\right)}{d c} & 0 & 0 \\
0 & \frac{P\left(U_{c}\right)}{c} & 0 \\
0 & 0 & \frac{P\left(U_{c}\right)}{c}
\end{array}\right) .
$$

Since

$$
n^{\leq 0}\left(\left.L_{c}\right|_{\operatorname{span}\left\{v_{1}, v_{2}, v_{3}\right\}}\right) \leq n^{-}\left(L_{c}\right) \leq 1
$$

by the index counting formula (2.13), so regardless of the sign of $\frac{d P\left(U_{c}\right)}{d c}$, we must have $P\left(U_{c}\right)>0$. The 3D analogue (see [60]) of the identity (11.78) is

$$
c P\left(U_{c}\right)=\int_{\mathbf{R}^{3}}\left|\frac{\partial U_{c}}{\partial x_{1}}\right|^{2} d x+\int_{\mathbf{R}^{3}} V\left(\left|U_{c}\right|\right)^{2} d x
$$

which is again not enough to conclude $P\left(U_{c}\right)>0$ when $V$ takes negative values.

## CHAPTER 12

## Appendix

In this appendix, we give some elementary properties of (2.1), which are mostly based on theoretical functional analysis arguments. They include some basic decomposition of the phase space, the well-posedness of (2.1), and the standard complexification procedure.

We start with some elementary properties of $L$. First we prove that $n^{-}(L)=$ $\operatorname{dim} X_{-}$in assumption (H2) is the maximal dimension of subspaces where $\langle L \cdot, \cdot\rangle<$ 0 .

Lemma 12.1. If $N \subset X$ is a subspace such that $\langle L u, u\rangle<0$ for all $u \in N \backslash\{0\}$, then $\operatorname{dim} N \leq n^{-}(L)$.

Proof. Let $X_{ \pm}$be given in (H2) and $P_{+, 0,-}$ be the projections associated to the decomposition $X=X_{+} \oplus \operatorname{ker} L \oplus X_{-}$. For any $u \in X, P_{-} u=0$ would imply $u \in \operatorname{ker} L \oplus X_{+}$and thus $\langle L u, u\rangle \geq 0$, so $u \notin N$. Therefore, $P_{-}: N \rightarrow X_{-}$is injective and in turn it implies $\operatorname{dim} N \leq \operatorname{dim} X_{-}$.

In order to proceed we have to introduce some notations. Given a closed subspace $Y \subset X$, let $i_{Y}: Y \rightarrow X$ be the embedding and then $i_{Y}^{*}: X^{*} \rightarrow Y^{*}$. Define

$$
\begin{align*}
& L_{Y}=i_{Y}^{*} L i_{Y}: Y \rightarrow Y^{*}, \\
& Y^{\perp_{L}}=\operatorname{ker}\left(i_{Y}^{*} L\right)=\left\{u \in X \mid\left\langle L u, i_{Y} v\right\rangle=\langle L u, v\rangle=0, \forall v \in Y\right\}, \tag{12.1}
\end{align*}
$$

which satisfy

$$
\begin{equation*}
L_{Y}^{*}=L_{Y} \text { and }\left\langle L_{Y} u, v\right\rangle=\langle L u, v\rangle, \forall u, v \in Y \tag{12.2}
\end{equation*}
$$

The following is a simple technical lemma.
Lemma 12.2. Assume (H1-3). Let $Y \subset X$ be a closed subspace.
(1) Suppose the quadratic form $\langle L \cdot, \cdot\rangle$ is non-degenerate (in the sense of (2.4)) on $Y$, then $X=Y \oplus Y^{\perp_{L}}$.
(2) Assume $\operatorname{dim} \operatorname{ker} L<\infty$ and $\operatorname{ker} L_{Y}=\{0\}$, then $\langle L \cdot, \cdot\rangle$ is non-degenerate on $Y$.
(3) If $X=\operatorname{ker} L \oplus Y$ then $\langle L \cdot, \cdot\rangle$ is non-degenerate on $Y$.

Proof. We first notice that $L_{Y}$ being an isomorphism implies $Y \cap Y^{\perp_{L}}=\{0\}$. For any $u \in X$, let

$$
u_{1}=L_{Y}^{-1} i_{Y}^{*} L u \in Y \Longrightarrow\left\langle L u_{1}-L u, v\right\rangle=0, \forall v \in Y
$$

and thus $u_{2}=u-u_{1} \in Y^{\perp_{L}}$ which implies $X=Y \oplus Y^{\perp_{L}}$.

In order prove the second statement, from the standard argument, it suffices to show that

$$
\begin{equation*}
\inf _{u \in Y \backslash\{0\}} \sup _{v \in Y \backslash\{0\}} \frac{|\langle L u, v\rangle|}{\|u\|\|v\|}>0 \tag{12.3}
\end{equation*}
$$

According to Remark 2.2 and the assumption of the lemma, there exist closed subspaces $X_{\leq 0}$ and $X_{+}$such that the decomposition $X=X_{\leq 0} \oplus X_{+}$is orthogonal with respect to both $(\cdot, \cdot)$ and $\langle L \cdot, \cdot\rangle, \operatorname{dim} X_{\leq 0}<\infty,\langle L u, u\rangle \leq 0$ for all $u \in X_{\leq 0}$, and for some $\delta>0,\langle L u, u\rangle \geq \delta\|u\|^{2}$ for all $u \in X_{+}$. This splitting is associated to the orthogonal projections $\mathcal{P}_{\leq 0,+}: X \rightarrow X_{\leq 0,+}$. Let $Y_{+}=Y \cap X_{+}$and

$$
Y_{1}=\left\{u \in Y \mid\langle L u, v\rangle=0, \forall v \in Y_{+}\right\}
$$

Clearly, $Y_{+}$and $Y_{1}$ are both closed subspaces of $Y$. Much as in the first statement, using the uniform positive definiteness of $\langle L u, u\rangle$ on $Y_{+}$, we have $Y=Y_{+} \oplus Y_{1}$ via

$$
u=u_{+}+\left(u-u_{+}\right), \text {where } u_{+}=L_{Y_{+}}^{-1} i_{Y_{+}}^{*} L u \in Y_{+}, \quad \forall u \in Y
$$

For any $u_{1} \in Y_{1} \backslash\{0\}$, let $x_{\leq 0,+}=\mathcal{P}_{X_{<0,+}} u_{1}$ and we have $u_{1}=x_{\leq 0}+x_{+}$. Since $\mathcal{P}_{X_{\leq 0}} u_{1}=0$ would imply $u_{1} \in \bar{X}_{+} \cap Y=Y_{+}$contradictory to $Y=Y_{+} \oplus Y_{1}$, we obtain that the linear mapping $\left.\mathcal{P}_{X_{<0}}\right|_{Y_{1}}$ is one-to-one. Therefore, $\operatorname{dim} Y_{1}<\infty$. From the definition of $Y_{1}$, if $u_{1} \in Y_{1} \backslash\{0\}$ satisfies that $\left\langle L u_{1}, v\right\rangle=0$ for all $v \in Y_{1}$, we would have $L_{Y} u_{1}=0$ which contradicts the assumption ker $L_{Y}=\{0\}$. Therefore, $\left.L_{Y}\right|_{Y_{1}}$ defines an isomorphism from $Y_{1}$ to $Y_{1}^{*}$ as $\operatorname{dim} Y_{1}<\infty$ and thus there exists $\delta^{\prime}>0$ such that for any $u_{1} \in Y_{1} \backslash\{0\}$, there exists $v \in Y_{1}$ such that $\left\langle L_{Y} u_{1}, v\right\rangle \geq \delta^{\prime}\left\|u_{1}\right\|\|v\|$.

Consider any $u=u_{1}+u_{+} \in Y$. If $\left\|u_{1}\right\| \geq\left\|u_{+}\right\|$, there exists $v \in Y_{1}$ such that

$$
\langle L u, v\rangle=\left\langle L u_{1}, v\right\rangle \geq \delta^{\prime}\left\|u_{1}\right\|\|v\| \geq \frac{\delta^{\prime}}{2}\|u\|\|v\| .
$$

If $\left\|u_{+}\right\| \geq\left\|u_{1}\right\|$, then let $v=u_{+}$and we have

$$
\langle L u, v\rangle=\left\langle L u_{+}, u_{+}\right\rangle \geq \delta\left\|u_{+}\right\|^{2} \geq \frac{\delta}{2}\|u\|\|v\| .
$$

Therefore, (12.3) is obtained and the second statement is proved.
Finally we prove the last statement. We first show the non-degeneracy of $\langle L \cdot, \cdot\rangle$ on $X_{+} \oplus X_{-}$though a standard procedure. The bounded symmetric quadratic form $\langle L \cdot, \cdot\rangle$ on $X_{+} \oplus X_{-}$induces bounded linear operators

$$
L_{\alpha, \beta}=i_{X_{\alpha}}^{*} L i_{X_{\beta}}: X_{\beta} \rightarrow X_{\alpha}^{*}, \quad \alpha, \beta \in\{+,-\}
$$

Since $L_{++}$and $-L_{--}$are both symmetric and bounded below, thus isomorphic, and $L_{+-}=L_{-+}^{*}$, so the same are true for

$$
L_{++}-L_{+-} L_{--}^{-1} L_{-+} \text {and }-\left(L_{--}-L_{-+} L_{++}^{-1} L_{+-}\right)
$$

It is easy to verify that

$$
L^{-1}=\left(L_{++}-L_{+-} L_{--}^{-1} L_{-+}\right)^{-1} i_{X_{+}}^{*}+\left(L_{--}-L_{-+} L_{++}^{-1} L_{+-}\right)^{-1} i_{X_{-}}^{*}
$$

is a bounded operator from $\left(X_{+} \oplus X_{-}\right)^{*}$ to $X_{+} \oplus X_{-}$. In general, if $X=\operatorname{ker} L \oplus Y$, there exists an isomorphism $T: X_{-} \oplus X_{+} \rightarrow \operatorname{ker} L$ such that $Y=\operatorname{graph}(T)$. The non-degeneracy of $\langle L \cdot, \cdot\rangle$ on $Y$ follows immediately from its non-degeneracy on $X_{-} \oplus X_{+}$. The proof of the lemma is complete.

Remark 12.1. The first statement in the lemma holds actually for any closed subspace $Y \subset X$ as long as $\langle L \cdot, \cdot\rangle$ is non-degenerate on $Y$. The finite dimensionality assumption on $\operatorname{ker} L$ is essential for the second statement in the above lemma. A counter example is

$$
X=l^{2} \oplus l^{2}, L=I \oplus 0, Y=\left\{\left(\left\{x_{n}\right\},\left\{y_{n}\right\}\right) \in X \left\lvert\, x_{n}=\frac{1}{n} y_{n}\right.\right\}
$$

for which $\operatorname{dim} \operatorname{ker} L=\infty, n^{-}(L)=0,\left.\operatorname{ker} L\right|_{Y}=\{0\}$, but $\langle L \cdot, \cdot\rangle$ is not nondegenerate on $Y$ in the sense of (2.4).

The next lemma will allow us to decompose equation (2.1).
Lemma 12.3. Suppose $X_{1,2} \subset X$ are closed subspaces satisfying $X=X_{1} \oplus X_{2}$. Let $P_{1,2}: X \rightarrow X_{1,2}$ be the associated projections, which imply $P_{1,2}^{*}: X_{1,2}^{*} \rightarrow X^{*}$, and
$J_{j k}=P_{j} J P_{k}^{*}: D\left(J_{j k}\right) \rightarrow X_{j}, D\left(J_{j k}\right)=\left(P_{k}^{*}\right)^{-1}\left(D(J) \cap P_{k}^{*} X_{k}^{*}\right), j, k=1,2$.
(1) If $\operatorname{ker} i_{X_{2}}^{*} \subset D(J)$, then $J_{11}$ and $J_{21}$ are bounded operators defined on $X_{1}^{*}$, $J_{11}^{*}=-J_{11}, J_{22}=-J_{22}^{*}$, and $J_{12}^{*}=-J_{21}$, and $J_{12}$ can be extended to the bounded operator $-J_{21}^{*}=J_{12}^{* *}$ defined on $X_{2}^{*}$.
(2) If $\left\langle L u_{1}, u_{2}\right\rangle=0$, for all $u_{j} \in X_{j}, j=1,2$, then $L X_{j} \subset \operatorname{ker} i_{X_{3-j}}^{*}, L_{X_{1,2}}$ satisfy (H2) on $X_{1,2}, n^{-}(L)=n^{-}\left(L_{X_{1}}\right)+n^{-}\left(L_{X_{2}}\right)$, and $\operatorname{ker} L=\operatorname{ker} L_{X_{1}} \oplus$ $\operatorname{ker} L_{X_{2}}$.
(3) Assume $\left\langle L u_{1}, u_{2}\right\rangle=0$, for all $u_{j} \in X_{j}, j=1,2$, and ker $i_{X_{2}}^{*} \subset D(J)$, then the combinations $\left(X_{j}, L_{X_{j}}, J_{j j}\right), j=1,2$, satisfy (H1-3).
Proof. For $j=1,2$, define $\tilde{X}_{j}^{*}$ as

$$
\begin{equation*}
\tilde{X}_{j}^{*}=P_{j}^{*} X_{j}^{*}=\operatorname{ker} i_{X_{3-j}}^{*}=\left\{f \in X^{*} \mid\langle f, u\rangle=0, \forall u \in X_{3-j}\right\} \subset X^{*} \tag{12.4}
\end{equation*}
$$

Clearly, it holds

$$
\begin{equation*}
i_{X_{1}} P_{1}+i_{X_{2}} P_{2}=I_{X}, \quad P_{1}^{*} i_{X_{1}}^{*}+P_{2}^{*} i_{X_{2}}^{*}=I_{X^{*}}, \quad X^{*}=\tilde{X}_{1}^{*} \oplus \tilde{X}_{2}^{*} \tag{12.5}
\end{equation*}
$$

Assume $\tilde{X}_{1}^{*}=P_{1}^{*} X_{1}^{*} \subset D(J)$. The Closed Graph Theorem implies that the closed operator $J P_{1}^{*}: X_{1}^{*} \rightarrow X$ is actually bounded, and thus $J_{11}$ and $J_{21}$ are bounded as well. The property $J_{11}^{*}=-J_{11}$ is obvious from $J^{*}=-J$ and the boundedness of $J_{11}$. We also obtain from this assumption and (12.5) that $D(J) \cap \tilde{X}_{2}^{*}$ is dense in $\tilde{X}_{2}^{*}$ and thus $J_{12}$ and $J_{22}$ are densely defined, as $P_{j}^{*}: X_{j}^{*} \rightarrow \tilde{X}_{j}^{*}$ is an isomorphism. It remains to prove $J_{12}^{*}=-J_{21}$ and $J_{22}^{*}=-J_{22}$.

Suppose $u=J_{12}^{*} g$, or equivalently, $g \in X_{1}^{*}$ and $u \in X_{2}$ satisfy, $\forall f \in D\left(J_{12}\right) \subset$ $X_{2}^{*}$,

$$
\begin{equation*}
\left\langle P_{2}^{*} f, i_{X_{2}} u-i_{X_{1}} P_{1} J P_{1}^{*} g\right\rangle=\langle f, u\rangle=\left\langle g, J_{12} f\right\rangle=\left\langle P_{1}^{*} g, J P_{2}^{*} f\right\rangle \tag{12.6}
\end{equation*}
$$

where we used $P_{j} i_{X_{j}}=i d$ and $P_{3-j} i_{X_{j}}=0$ on $X_{j}$. For any $h \in X_{1}^{*}$, we have

$$
\left\langle P_{1}^{*} h, i_{x_{2}} u-i_{X_{1}} P_{1} J P_{1}^{*} g\right\rangle=\left\langle P_{1}^{*} g, J P_{1}^{*} h\right\rangle
$$

Therefore, (12.5) and (12.6) imply $u=J_{12}^{*} g$ is equivalent to

$$
\begin{aligned}
& \left\langle\gamma, i_{X_{2}} u-i_{X_{1}} P_{1} J P_{1}^{*} g\right\rangle=\left\langle P_{1}^{*} g, J \gamma\right\rangle, \quad \forall \gamma \in D(J) \\
\Longleftrightarrow & i_{X_{2}} u-i_{X_{1}} P_{1} J P_{1}^{*} g=J^{*} P_{1}^{*} g=-J P_{1}^{*} g \\
\Longleftrightarrow & u=-P_{2} J P_{1}^{*} g=-J_{21} g
\end{aligned}
$$

Therefore, $J_{12}^{*}=-J_{21}$.

Similarly, using the assumption $\tilde{X}_{1}^{*} \subset D(J)$, one can prove $u=J_{22}^{*} g \in X_{2}$, $g \in X_{2}^{*}$, if and only if

$$
i_{X_{2}} u+i_{X_{1}} J_{21}^{*} g=J^{*} P_{2}^{*} g \Longleftrightarrow u=-P_{2} J P_{2}^{*} g=-J_{22} g .
$$

Therefore, we obtain $J_{22}^{*}=-J_{22}$.
Assume $\left\langle L u_{1}, u_{2}\right\rangle=0$, for all $u_{j} \in X_{j}, j=1,2$. As a direct consequence, we have $L X_{j} \subset \tilde{X}_{j}^{*}$, which, along with (12.5), immediately implies

$$
L=P_{1}^{*} L_{X_{1}} P_{1}+P_{2}^{*} L_{X_{2}} P_{2}, \quad P_{j}^{*} L_{X_{j}} P_{j}(X) \subset \tilde{X}_{j}
$$

which in turn yield

$$
\operatorname{ker} L=\operatorname{ker} L_{X_{1}} \oplus \operatorname{ker} L_{X_{2}}, \quad \operatorname{ker} L_{X_{j}}=X_{j} \cap \operatorname{ker} L, j=1,2
$$

Let

$$
Y_{1,2}=\left\{u \in X_{1,2} \mid(u, v)=0, \forall v \in \operatorname{ker} L_{X_{1,2}}\right\}, Y=Y_{1} \oplus Y_{2}
$$

which implies

$$
X=Y \oplus \operatorname{ker} L=Y_{1} \oplus Y_{2} \oplus \operatorname{ker} L
$$

and

$$
\left\langle L y_{j}, y_{1}^{\prime}+y_{2}^{\prime}+u\right\rangle=\left\langle L y_{j}, y_{j}^{\prime}\right\rangle=\left\langle L_{Y_{j}} y_{j}, y_{j}^{\prime}\right\rangle
$$

for any $y_{j}, y_{j}^{\prime} \in Y_{j}, j=1,2$, and $u \in \operatorname{ker} L$. Let $P_{Y_{1,2,0}}$ be the projections associated to this decomposition, then we have

$$
L\left(Y_{j}\right) \subset \tilde{Y}_{j}^{*} \triangleq P_{Y_{j}}^{*} Y_{j}^{*}=\operatorname{ker} i_{\operatorname{ker} L \oplus Y_{3-j}}^{*}, \quad j=1,2
$$

Assumption (H2) implies that $\left.L\right|_{Y}: Y \rightarrow R(L)$ is an isomorphism to the closed subspace $R(L) \subset X^{*}$. Therefore, $L\left(Y_{1,2}\right) \subset \tilde{Y}_{1,2}^{*}$ are closed subspaces and $\left.L\right|_{Y_{1,2}}$ : $Y_{1,2} \rightarrow L\left(Y_{1,2}\right)$ are isomorphisms. It implies that $L_{Y_{1,2}}$ are isomorphisms from $Y_{1,2}$ to closed subspaces $L_{Y_{1,2}}\left(Y_{1,2}\right) \subset Y_{1,2}^{*}$. Due to their boundedness and symmetry, we obtain that $L_{Y_{1,2}} Y_{1,2}$ is equal to the orthogonal complement of $\operatorname{ker} L_{Y_{1,2}}^{*}=$ $\operatorname{ker} L_{Y_{1,2}}=\{0\}$. So $L_{Y_{1,2}}: Y_{1,2} \rightarrow Y_{1,2}^{*}$ are isomorphisms, which induce bounded non-degenerate symmetric quadratic forms on $Y_{1,2}$. From the standard theory on symmetric quadratic forms, $Y_{j}, j=1,2$, can be split into $Y_{j}=Y_{j+} \oplus Y_{j-}$, where closed subspaces $Y_{j \pm}$ are orthogonal with respect to both $(\cdot, \cdot)$ and $\langle L \cdot, \cdot\rangle$. Moreover, there exists $\delta>0$ such that

$$
\pm\left\langle L_{X_{j}} u, u\right\rangle= \pm\langle L u, u\rangle \geq \delta\|u\|^{2}, \forall u \in Y_{j \pm}
$$

This proves that $X_{j}$ satisfies (H2) with

$$
X_{j}=Y_{J-} \oplus \operatorname{ker} L_{X_{j}} \oplus Y_{j+}, \quad j=1,2
$$

Finally, since $X=X_{1} \oplus X_{2}$, there exists $C>0$ such that,

$$
\left\|u_{1}\right\|^{2}+\left\|u_{2}\right\|^{2} \leq C\left\|u_{1}+u_{2}\right\|^{2}, \quad \forall u_{1,2} \in X_{1,2}
$$

Therefore, the splitting

$$
X=\left(Y_{1-} \oplus Y_{2-}\right) \oplus \operatorname{ker} L \oplus\left(Y_{1+} \oplus Y_{2+}\right)
$$

satisfies the properties in (H2), which implies $n^{-}(L)=n^{-}\left(L_{X_{1}}\right)+n^{-}\left(L_{X_{2}}\right)$.
Finally, assume $\left\langle L u_{1}, u_{2}\right\rangle=0$, for all $u_{j} \in X_{j}, j=1,2$, and $P_{1}^{*} X_{1}^{*} \subset D(J)$. To complete the proof of the lemma, we only need to show that (H3) is satisfied by $\left(X_{j}, L_{X_{j}}, J_{j j}\right), j=1,2$. This is obvious for $j=1$, as $J_{11}$ is a bounded operator, and thus we only need to work on $j=2$. Let $X_{ \pm} \subset X$ be the closed subspaces assumed in (H2-3) and $Z=X_{-} \oplus X_{+}$. Since $X=\operatorname{ker} L \oplus Z=\operatorname{ker} L \oplus Y, Z$ can be represented as the graph of a bounded linear operator from $Y$ to ker $L$.

As $\operatorname{ker} L=\operatorname{ker} L_{X_{1}} \oplus \operatorname{ker} L_{X_{2}}$ and $Y=Y_{1} \oplus Y_{2}$, there exist bounded operators $S_{j k}: Y_{k} \rightarrow \operatorname{ker} L_{X_{j}}$ such that

$$
Z=\left\{y_{1}+y_{2}+\Sigma_{j, k=1}^{2} S_{j k} y_{k} \mid y_{1,2} \in Y_{1,2}\right\}
$$

We will first show

$$
\begin{equation*}
W \triangleq\left\{f \in X_{2}^{*} \mid\langle f, u\rangle=0, u \in Z_{2}\right\} \subset D\left(J_{22}\right) \tag{12.7}
\end{equation*}
$$

where

$$
Z_{2}=\left\{y_{2}+S_{22} y_{2} \mid y_{2} \in Y_{2}\right\} \subset X_{2} .
$$

Trivially extend $S_{j k}$ to be an operator from $X_{k}$ to $\operatorname{ker} L_{X_{j}} \subset X_{j}$ via

$$
S_{j k}\left(y_{k}+v_{k}\right)=S_{j k} y_{k}, \quad \forall y_{k} \in Y_{k}, \quad v_{k} \in \operatorname{ker} L_{X_{k}} .
$$

It leads to $S_{j k} S_{k l}=0, \forall j, k, l=1,2$. Given any $f \in W \subset X_{2}^{*}$, one may compute, for any

$$
u=y_{1}+y_{2}+\Sigma_{j, k=1}^{2} S_{j k} y_{k} \in Z
$$

using the definition of $W$, and the property of the extensions of $S_{j k}$,

$$
\begin{aligned}
\left\langle P_{2}^{*} f-P_{1}^{*} S_{21}^{*} f, u\right\rangle & =\left\langle f, y_{2}+S_{21} y_{1}+S_{22} y_{2}\right\rangle-\left\langle S_{21}^{*} f, y_{1}+S_{11} y_{1}+S_{12} y_{2}\right\rangle \\
& =\left\langle f, S_{21} y_{1}\right\rangle-\left\langle f, S_{21} y_{1}+S_{21} S_{11} y_{1}+S_{21} S_{12} y_{2}\right\rangle=0 .
\end{aligned}
$$

Therefore, (H3) implies $P_{2}^{*} f-P_{1}^{*} S_{21}^{*} f \in D(J)$. Since we assume $P_{1}^{*} X_{1}^{*} \subset D(J)$, we obtain $P_{2}^{*} f \in D(J)$ and thus $f \in D\left(J_{22}\right)$ which proves (12.7).

Since $y_{2} \rightarrow y_{2}+S_{22} y_{2}$ is an isomorphism from $Y_{2}$ to $Z_{2}$,

$$
\left\langle L\left(y_{2}+S_{22} y_{2}\right), y_{2}^{\prime}+S_{22} y_{2}^{\prime}\right\rangle=\left\langle L y_{2}, y_{2}^{\prime}\right\rangle
$$

and $L_{Y_{2}}$ is isomorphic, we have $\langle L \cdot, \cdot\rangle$ is non-degenerate on $Z_{2}$ and $L_{Z_{2}}$ is also an isomorphism. Therefore, there exist closed subspaces $X_{2 \pm} \subset Z_{2}$ and $\delta>0$ such that $Z_{2}=X_{2-} \oplus X_{2+}, \operatorname{dim} X_{2-}=n^{-}\left(L_{X_{2}}\right)$, and $\pm\left\langle L_{X_{2}} u, u\right\rangle \geq \delta\|u\|^{2}$, for any $u \in X_{2 \pm}$. It along with (12.7) and $X_{2}=Z_{2} \oplus \operatorname{ker} L_{X_{2}}$ completes the proof of the lemma.

Remark 12.2. Under assumptions $\left\langle L u_{1}, u_{2}\right\rangle=0$, for all $u_{j} \in X_{j}, j=1,2$, and $P_{1}^{*} X_{1}^{*} \subset D(J),\left(X_{j}, L_{X_{j}}, J_{j j}\right), j=1,2$, satisfies the same hypothesis (H1-H3) as $(X, L, J)$ and $n^{-}(L)=n^{-}\left(L_{X_{1}}\right)+n^{-}\left(L_{X_{2}}\right)$. Moreover, it is easily verified based on these assumptions that $J_{j j} L_{X_{j}}=\left.P_{j} J L\right|_{X_{j}}$. Therefore, this lemma would often be applied to reduce the problem to subspaces when $J L\left(X_{1}\right) \subset X_{1}$, which implies JL has certain upper triangular structure.

Corollary 12.1. LJ: $D(J) \rightarrow X^{*}$ is a closed operator and consequently $(J L)^{*}=-L J$.

Proof. Let $X_{ \pm}$and ker $L$ satisfy the requirements in (H2-3) and let $X_{1}=$ ker $L$ and $X_{2}=X_{-} \oplus X_{+}$. Clearly, we have, $L_{X_{1}}=0,\left\langle L u_{1}, u_{2}\right\rangle=0$, for all $u_{j} \in X_{j}$, $j=1,2$, and $P_{1}^{*} X_{1}^{*} \subset D(J)$ due to (H3). Using

$$
i_{X_{1}} P_{1}+i_{X_{2}} P_{2}=I_{X}, \quad L i_{X_{1}}=0, \quad i_{X_{1}}^{*} L=0
$$

$L J$ can be rewritten in this decomposition

$$
L J \gamma=P_{2}^{*} L_{X_{2}} J_{21} i_{X_{1}}^{*} \gamma+P_{2}^{*} L_{X_{2}} J_{22} i_{X_{2}}^{*} \gamma, \quad \forall \gamma \in X^{*}
$$

which is equivalent to using the blockwise decomposition of $J$ and $L$. Since $J_{21}$ is continuous, $P_{2}^{*} L_{X_{2}} J_{21} i_{X_{1}}^{*}$ is continuous too. Moreover, the facts that $L_{X_{2}}: X_{2} \rightarrow$ $X_{2}^{*}$ is an isomorphism, $P_{2}^{*}$ has a continuous left inverse $i_{X_{2}}^{*}$ as $P_{2} i_{X_{2}}=I_{X_{2}}$, along
with the closedness of $J_{22}$ imply that $P_{2}^{*} L_{X_{2}} J_{22}$ and thus $P_{2}^{*} L_{X_{2}} J_{22} i_{X_{2}}^{*}$ is a closed operator. Therefore, $L J$ is closed.

Since $(L J)^{*}=J L$ is densely defined and thus $(L J)^{* *}=-(J L)^{*}$ is well defined. The closeness of $L J$ implies $L J=(L J)^{* *}=-(J L)^{*}$.

Remark 12.3. We would like to point out that, in the proof Lemma 12.3 and Corollary 12.1, we do not use the assumption that $n^{-}(L)<\infty$. Therefore, they actually hold even if $n^{-}(L)=\infty$ except that $n^{-}\left(L_{X_{1,2}}\right)$ might be $\infty$.

The following is a simple, but useful, technical lemma.
Lemma 12.4. There exist closed subspaces $X_{ \pm} \subset X$ satisfying the properties in (H2-3) and in addition,
(1) $X=X_{0} \oplus X_{-} \oplus X_{+}$is a L-orthogonal splitting with associated projections $P_{0, \pm}$, where $X_{0}=\operatorname{ker} L$;
(2) $L_{X_{ \pm}}: X_{ \pm} \rightarrow X_{ \pm}^{*}$ are isomorphic; and
(3) $\tilde{X}_{0,-}^{*} \subset D(J)$ and $D(J) \cap \tilde{X}_{+}^{*}$ is dense in $\tilde{X}_{+}^{*}$, where $\tilde{X}_{ \pm, 0}^{*} \triangleq P_{ \pm, 0}^{*} X_{ \pm, 0}^{*}$ (see (12.1) and (12.4)).
Proof. Let $Y_{ \pm} \subset X$ be closed subspaces satisfying hypothesis (H2-3). Let $Y=Y_{-} \oplus Y_{+}, P: X \rightarrow Y$ be the projection associated to the decomposition $X=X_{0} \oplus Y, \tilde{X}_{0}^{*}=(I-P)^{*} X_{0}^{*}$, and $\tilde{Y}^{*}=P^{*} Y^{*}$, which are closed subspaces. According to (H3), we have $\tilde{X}_{0}^{*} \subset D(J)$. Consequently, $\tilde{Y}^{*} \cap D(J)$ is dense in $\tilde{Y}^{*}$ as $X^{*}=\tilde{X}_{0}^{*} \oplus \tilde{Y}^{*}$. Our assumptions imply $L_{Y}: Y \rightarrow \tilde{Y}^{*}$ is an isomorphism, which induces a bounded symmetric quadratic form on $Y$ with Morse index equal to $n^{-}(L)$. Therefore, there exists a closed subspace $X_{-} \subset Y$ such that $\operatorname{dim} X_{-}=$ $n^{-}(L), L\left(X_{-}\right) \subset D(J)$, and $\langle L u, u\rangle \leq-\delta\|u\|^{2}$, for all $u \in X_{-}$. Let

$$
X_{+}=\left\{u \in Y \mid\langle L u, v\rangle=0, \forall v \in X_{-}\right\}
$$

Since $L$ is uniformly negative on $X_{-}$, Lemma 12.2 implies the $L$-orthogonal splitting $Y=X_{-} \oplus X_{+}$and thus the $L$-orthogonal decomposition $X=X_{0} \oplus X_{-} \oplus X_{+}$as well. The rest of the proof follows easily from the facts that $L_{Y}$ is isomorphic, $X^{*}=\tilde{X}_{0}^{*} \oplus \tilde{X}_{-}^{*} \oplus \tilde{X}_{+}^{*}, \operatorname{dim} X_{-}=n^{-}(L), \tilde{X}_{0}^{*} \subset D(J)$, and $\tilde{X}_{-}^{*}=L\left(X_{-}\right) \subset D(J)$.

Remark 12.4. Under assumption (2.2), it is possible to choose $X_{ \pm}$such that $X_{+} \oplus X_{-}=(\operatorname{ker} L)^{\perp}$ satisfies all properties in Lemma 12.4, where $(\operatorname{ker} L)^{\perp}$ is defined in (2.3). In fact, let $Y=(\operatorname{ker} L)^{\perp}$, then (2.2) implies that the splitting $X=\operatorname{ker} L \oplus Y$ satisfies all assumptions in Lemma 12.3. The rest of the construction of $X_{ \pm} \subset Y=(\operatorname{ker} L)^{\perp}$ follows in exactly the same procedure as in the proof of Lemma 12.4.

In order to establish the well-posedness of the linear equation in the next, we start with the following lemma.

Lemma 12.5. There exists an equivalent inner product $(\cdot, \cdot)_{L}$ on $X$, a linear operator $A: D(J L) \rightarrow X$ which is anti-self-adjoint with respect to $(\cdot, \cdot)_{L}$, and a bound linear operator $B: X \rightarrow X$ such that $J L=A+B$.

Proof. Let $X=X_{-} \oplus X_{0} \oplus X_{+}$be a decomposition as given in Lemma 12.4 with $X_{0}=\operatorname{ker} L$. Let

$$
L_{ \pm}= \pm P_{ \pm}^{*} i_{X_{ \pm}}^{*} L i_{X_{ \pm}} P_{ \pm}: X \rightarrow X^{*}
$$

which satisfy

$$
L_{ \pm}^{*}=L_{ \pm}, \quad L=L_{+}-L_{-}, \quad\left\langle L_{ \pm} u, v\right\rangle= \pm\left\langle L_{X_{ \pm}} u, v\right\rangle= \pm\langle L u, v\rangle, \forall u, v \in X_{ \pm}
$$

Let $R: X \rightarrow X^{*}$ be the isomorphism corresponding to $(\cdot, \cdot)$ through the Riesz Representation Theorem and

$$
L_{0}=P_{0}^{*} i_{X_{0}}^{*} R i_{X_{0}} P_{X_{0}}: X \rightarrow X^{*} \longleftrightarrow\left\langle L_{0} u, v\right\rangle=\left(P_{0} u, P_{0} v\right)
$$

From Lemma 12.4 and assumptions (H2-3), it is easy to verify that

$$
(u, v)_{L} \triangleq\left\langle\left(L_{+}+L_{-}+L_{0}\right) u, v\right\rangle=\left\langle L_{+} u, v\right\rangle+\left\langle L_{-} u, v\right\rangle+\left(P_{0} u, P_{0} v\right)
$$

is uniformly positive and defines an equivalent inner product on $X$. Let

$$
A=J\left(L_{+}+L_{-}+L_{0}\right)=J L+2 J L_{-}+J L_{0} \triangleq J L-B
$$

Since $P_{0,-}^{*} X_{0,-}^{*} \subset D(J)$, the Closed Graph Theorem implies that $B$ is bounded. If $\operatorname{dim} \operatorname{ker} L<\infty, B$ is obviously of finite rank. The proof of the lemma is complete.

A direct consequence of this lemma is the well-posedness of equation (2.1) which follows from the standard perturbation theory of semigroups.

Proposition 12.1. JL generates a $C^{0}$ group $e^{t J L}$ of bounded linear operators on $X$.

Complexification. For considerations where complex eigenvalues are involved, we have to work with the standard complexification of $X$ and the associated operators. Let

$$
\tilde{X}=\left\{x=x_{1}+i x_{2} \mid x_{1,2} \in X\right\} \text { with } \overline{x_{1}+i x_{2}}=x_{1}-i x_{2}
$$

equipped with the complexified inner product

$$
\left(x_{1}+i x_{2}, x_{1}^{\prime}+i x_{2}^{\prime}\right)=\left(x_{1}, x_{1}^{\prime}\right)+\left(x_{2}, x_{2}^{\prime}\right)+i\left(\left(x_{2}, x_{1}^{\prime}\right)-\left(x_{1}, x_{2}^{\prime}\right)\right)
$$

Instead of complexifying $L$ as a linear operator directly, it is much more convenient for us to complexify its corresponding real symmetric quadratic form $\langle L u, v\rangle$ into a complex Hermitian symmetric form

$$
\begin{align*}
& \mathcal{B}\left(x_{1}^{\prime}+i x_{2}^{\prime}, x_{1}+i x_{2}\right)=\left\langle\tilde{L}\left(x_{1}+i x_{2}\right),\left(x_{1}^{\prime}+i x_{2}^{\prime}\right)\right\rangle \\
= & \left\langle L x_{1}, x_{1}^{\prime}\right\rangle+\left\langle L x_{2}, x_{2}^{\prime}\right\rangle+i\left(\left\langle L x_{1}, x_{2}^{\prime}\right\rangle-\left\langle L x_{2}, x_{1}^{\prime}\right\rangle\right) \tag{12.8}
\end{align*}
$$

for any $x_{1,2}, x_{1,2}^{\prime} \in X$. Accordingly $L$ is complexified to a (anti-linear) mapping $\tilde{L}$ from $\tilde{X}$ to $\tilde{X}^{*}$ satisfying

$$
\begin{equation*}
\tilde{L}\left(c x+c^{\prime} x^{\prime}\right)=\bar{c} \tilde{L} x+\bar{c}^{\prime} \tilde{L} x^{\prime} \tag{12.9}
\end{equation*}
$$

A similar complexification can also be carried out for $J$ corresponding to a Hermitian symmetric form on $\tilde{X}^{*}$ and a (anti-linear) mapping from $\tilde{X}^{*} \rightarrow \tilde{X}$.

The composition $\tilde{J} \circ \tilde{L}$ of (anti-linear) mappings $\tilde{J}$ and $\tilde{L}$ is a closed complex linear operator from $D(\tilde{J} \tilde{L}) \subset \tilde{X}$ to $\tilde{X}$. The fact that $\widetilde{J L}$ is anti-symmetric with respect to the Hermitian symmetric form $\langle\tilde{L} u, v\rangle$, that is,

$$
\begin{equation*}
\langle\tilde{L}(\tilde{J} \tilde{L} u), v\rangle=-\langle\tilde{L} u, \tilde{J} \tilde{L} v\rangle \tag{12.10}
\end{equation*}
$$

will be used frequently. According to Corollary 12.1, the dual operator of $\tilde{J} \tilde{L}$ is given by

$$
\begin{equation*}
(\tilde{J} \tilde{L})^{*}=-\tilde{L} \tilde{J} \tag{12.11}
\end{equation*}
$$

It is easy to verify that $\tilde{L}, \tilde{J}, \tilde{J} \tilde{L}$ and $\tilde{L} \tilde{J}$ are real in the sense

$$
\begin{equation*}
\overline{\left\langle\tilde{L} x, x^{\prime}\right\rangle}=\left\langle\tilde{L} \bar{x}, \overline{x^{\prime}}\right\rangle, \overline{\langle f, \tilde{J} g\rangle}=\langle\bar{f}, \tilde{J} \bar{g}\rangle, \overline{\tilde{J} \tilde{L} x}=\tilde{J} \tilde{L} \bar{x}, \overline{\tilde{L}} \tilde{J} x=\tilde{L} \tilde{J} \bar{x} . \tag{12.12}
\end{equation*}
$$

This implies that the spectrum of $\tilde{J} \tilde{L}$ and $\tilde{L} \tilde{J}$ are symmetric about the real axis in the complex plane.

Remark 12.5. In fact, on the complexified Hilbert space $\tilde{X}$ (or on $\tilde{X}^{*}$ ), a linear operator or a Hermitian form is the complexification of a (real) operator or a symmetric quadratic form on $X$ (or on $\tilde{X}^{*}$ ) if and only if (12.12) holds.

In the rest of the paper, with slight abuse of notations, we will write $X, J L,\langle L u, v\rangle$ also for their complexifications unless confusion might occur.

Remark 12.6. The linear group of bounded operators e ${ }^{t J L}$ obtained in Proposition 12.1 is also complexified accordingly when needed.

Remark 12.7. Exactly the same statements in Lemma 12.2, 3.1, 3.2 hold in the complexified framework.

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