# Seifert Conjecture in the Even Convex Case 

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#### Abstract

In this paper, we prove that there exist at least $n$ geometrically distinct brake orbits on every $C^{2}$ compact convex symmetric hypersurface $\Sigma$ in $\mathbb{R}^{2 n}$ satisfying the reversible condition $N \Sigma=\Sigma$ with $N=\operatorname{diag}\left(-I_{n}, I_{n}\right)$. As a consequence, we show that if the Hamiltonian function is convex and even, then Seifert conjecture of 1948 on the multiplicity of brake orbits holds for any positive integer $n$. © 2014 Wiley Periodicals, Inc.


## 1 Introduction

For the standard symplectic space $\left(\mathbb{R}^{2 n}, \omega_{0}\right)$ with $\omega_{0}(x, y)=\langle J x, y\rangle$, where $J=\left(\begin{array}{cc}0 & -I \\ I & 0\end{array}\right)$ is the standard symplectic matrix and $I$ is the $n \times n$ identity matrix, an involution matrix defined by $N=\left(\begin{array}{rr}-I & 0 \\ 0 & I\end{array}\right)$ is clearly antisymplectic, i.e., $N J=$ $-J N$. The fixed point set of $N$ and $-N$ are the Lagrangian subspaces $L_{0}=$ $\{0\} \times \mathbb{R}^{n}$ and $L_{1}=\mathbb{R}^{n} \times\{0\}$ of $\left(\mathbb{R}^{2 n}, \omega_{0}\right)$, respectively.

Suppose $H \in C^{2}\left(\mathbb{R}^{2 n} \backslash\{0\}, \mathbb{R}\right) \cap C^{1}\left(\mathbb{R}^{2 n}, \mathbb{R}\right)$ satisfies the reversible condition

$$
\begin{equation*}
H(N x)=H(x) \quad \forall x \in \mathbb{R}^{2 n} \tag{1.1}
\end{equation*}
$$

We consider the following fixed energy problem of a nonlinear Hamiltonian system with Lagrangian boundary conditions:

$$
\begin{align*}
\dot{x}(t) & =J H^{\prime}(x(t)),  \tag{1.2}\\
H(x(t)) & =h,  \tag{1.3}\\
x(0) & \in L_{0}, x(\tau / 2) \in L_{0} . \tag{1.4}
\end{align*}
$$

It is clear that a solution $(\tau, x)$ of (1.2)-(1.4) is a characteristic chord on the contact submanifold $\Sigma:=H^{-1}(h)=\left\{y \in \mathbb{R}^{2 n} \mid H(y)=h\right\}$ of $\left(\mathbb{R}^{2 n}, \omega_{0}\right)$ and satisfies

$$
\begin{align*}
x(-t) & =N x(t),  \tag{1.5}\\
x(\tau+t) & =x(t) . \tag{1.6}
\end{align*}
$$

In this paper this kind of $\tau$-periodic characteristic $(\tau, x)$ is called a brake orbit on the hypersurface $\Sigma$. We denote by $\mathcal{J}_{b}(\Sigma, H)$ the set of all brake orbits on $\Sigma$. Two brake orbits $\left(\tau_{i}, x_{i}\right) \in \mathcal{J}_{b}(\Sigma, H), i=1,2$, are equivalent if the two brake
orbits are geometrically the same, i.e., $x_{1}(\mathbb{R})=x_{2}(\mathbb{R})$. We denote by $[(\tau, x)]$ the equivalence class of $(\tau, x) \in \mathcal{J}_{b}(\Sigma, H)$ in this equivalence relation and by $\widetilde{\mathcal{J}}_{b}(\Sigma, H)$ the set of $[(\tau, x)]$ for all $(\tau, x) \in \mathcal{J}_{b}(\Sigma, H)$. In fact, $\widetilde{\mathcal{J}}_{b}(\Sigma, H)$ is the set of geometrically distinct brake orbits on $\Sigma$, which is independent of the choice of $H$. So from now on we simply denote it by $\widetilde{\mathcal{J}}_{b}(\Sigma)$ and in the notation $[(\tau, x)]$ we always assume $x$ has minimal period $\tau$. We also denote by $\widetilde{\mathcal{J}}(\Sigma)$ the set of all geometrically distinct closed characteristics on $\Sigma$. The number of elements in a set $S$ is denoted by ${ }^{\#} S$. It is well-known that ${ }^{\#} \tilde{\mathcal{J}}_{b}(\Sigma)$ (and also ${ }^{\#} \tilde{\mathcal{J}}(\Sigma)$ ) is only dependent on $\Sigma$; that is to say, for simplicity we take $h=1$ if $H$ and $G$ are two $C^{2}$-functions satisfying (1.1) and $\Sigma_{H}:=H^{-1}(1)=\Sigma_{G}:=G^{-1}(1)$; then ${ }^{\#} \mathcal{J}_{b}\left(\Sigma_{H}\right)={ }^{\#} \mathcal{J}_{b}\left(\Sigma_{G}\right)$.

So we can consider the brake orbit problem in a more general setting. Let $\Sigma$ be a $C^{2}$ compact hypersurface in $\mathbb{R}^{2 n}$ bounding a compact set $C$ with nonempty interior. Suppose $\Sigma$ has nonvanishing Gaussian curvature and satisfies the reversible condition $N\left(\Sigma-x_{0}\right)=\Sigma-x_{0}:=\left\{x-x_{0} \mid x \in \Sigma\right\}$ for some $x_{0} \in C$. Without loss of generality, we may assume $x_{0}=0$. We denote the set of all such hypersurfaces in $\mathbb{R}^{2 n}$ by $\mathcal{H}_{b}(2 n)$. For $x \in \Sigma$, let $n_{\Sigma}(x)$ be the unit outward normal vector at $x \in \Sigma$. Note that here by the reversible condition there holds $n_{\Sigma}(N x)=N n_{\Sigma}(x)$. We consider the dynamics problem of finding $\tau>0$ and a $C^{1}$ smooth curve $x:[0, \tau] \rightarrow \mathbb{R}^{2 n}$ such that

$$
\begin{gather*}
\dot{x}(t)=\operatorname{Jn}_{\Sigma}(x(t)), \quad x(t) \in \Sigma,  \tag{1.7}\\
x(-t)=N x(t), \quad x(\tau+t)=x(t), \quad \text { for all } t \in \mathbb{R} . \tag{1.8}
\end{gather*}
$$

A solution $(\tau, x)$ of the problem (1.7)-(1.8) determines a brake orbit on $\Sigma$.
Definition 1.1. We denote by

$$
\begin{align*}
\mathcal{H}_{b}^{c}(2 n) & =\left\{\Sigma \in \mathcal{H}_{b}(2 n) \mid \Sigma \text { is strictly convex }\right\},  \tag{1.9}\\
\mathcal{H}_{b}^{s, c}(2 n) & =\left\{\Sigma \in \mathcal{H}_{b}^{c}(2 n) \mid-\Sigma=\Sigma\right\} \tag{1.10}
\end{align*}
$$

The main result of this paper is the following:
Theorem 1.2. For any $\Sigma \in \mathcal{H}_{b}^{s, c}(2 n)$ there holds

$$
\# \tilde{\mathcal{J}}_{b}(\Sigma) \geq n .
$$

Remark 1.3. Theorem 1.2 is a kind of multiplicity result related to the Arnold chord conjecture. The Arnold chord conjecture is an existence result that was proved by K. Mohnke in [24]. Another kind of multiplicity result related to the Arnold chord conjecture was proved in [11].

### 1.1 Seifert Conjecture

Let us recall the famous conjecture proposed by H. Seifert in his pioneer work [26] concerning the multiplicity of brake orbits in certain Hamiltonian systems in $\mathbb{R}^{2 n}$.

As a special case of (1.1), we assume $H \in C^{2}\left(\mathbb{R}^{2 n}, \mathbb{R}\right)$ possesses the form

$$
\begin{equation*}
H(p, q)=\frac{1}{2} A(q) p \cdot p+V(q) \tag{1.11}
\end{equation*}
$$

where $p, q \in \mathbb{R}^{n}, A(q)$ is a positive definite $n \times n$ for any $q \in \mathbb{R}^{n}$, and $A$ is $C^{2}$, and $V \in C^{2}\left(\mathbb{R}^{n}, \mathbb{R}\right)$ is the potential energy. It is clear that a solution of the Hamiltonian system

$$
\begin{gather*}
\dot{x}=J H^{\prime}(x), \quad x=(p, q),  \tag{1.12}\\
p(0)=p\left(\frac{\tau}{2}\right)=0, \tag{1.13}
\end{gather*}
$$

is a brake orbit. Moreover, if $h$ is the total energy of a brake orbit $(q, p)$, i.e., $H(p(t), q(t))=h$ and $V(q(0))=V(q(\tau))=h$, then $q(t) \in \bar{\Omega} \equiv\left\{q \in \mathbb{R}^{n} \mid\right.$ $V(q) \leq h\}$ for all $t \in \mathbb{R}$.

In [26] of 1948, H. Seifert studied the existence of brake orbit for system (1.12)(1.13) with the Hamiltonian function $H$ in the form of (1.11) and proved that $\overline{\mathcal{J}_{b}}(\Sigma) \neq \varnothing$ provided $V^{\prime} \neq 0$ on $\partial \Omega, V$ is analytic and $\bar{\Omega}$ is bounded and homeomorphic to the unit ball $B_{1}^{n}(0)$ in $\mathbb{R}^{n}$. Then in the same paper he proposed the following conjecture which is still open for $n \geq 2$ now:

$$
\# \tilde{\mathcal{J}}_{b}(\Sigma) \geq n \text { under the same conditions. }
$$

We note that for the Hamiltonian function

$$
H(p, q)=\frac{1}{2}|p|^{2}+\sum_{j=1}^{n} a_{j}^{2} q_{j}^{2}, \quad q, p \in \mathbb{R}^{n},
$$

where $a_{i} / a_{j} \notin \mathbb{Q}$ for all $i \neq j$ and $q=\left(q_{1}, q_{2}, \ldots, q_{n}\right)$. There are exactly $n$ geometrically distinct brake orbits on the energy hypersurface $\Sigma=H^{-1}(h)$.

### 1.2 Some Related Results since 1948

As a special case, letting $A(q)=I$ in 1.11 , the problem corresponds to the following classical fixed energy problem of the second-order autonomous Hamiltonian system

$$
\begin{gather*}
\ddot{q}(t)+V^{\prime}(q(t))=0 \quad \text { for } q(t) \in \Omega,  \tag{1.14}\\
\frac{1}{2}|\dot{q}(t)|^{2}+V(q(t))=h \quad \forall t \in \mathbb{R},  \tag{1.15}\\
\dot{q}(0)=\dot{q}\left(\frac{\tau}{2}\right)=0, \tag{1.16}
\end{gather*}
$$

where $V \in C^{2}\left(\mathbb{R}^{n}, \mathbb{R}\right)$ and $h$ is constant such that $\Omega \equiv\left\{q \in \mathbb{R}^{n} \mid V(q)<h\right\}$ is nonempty, bounded, and connected.

A solution $(\tau, q)$ of (1.14)-(1.16) is still called a brake orbit in $\bar{\Omega}$. Two brake orbits $q_{1}$ and $q_{2}: \mathbb{R} \rightarrow \mathbb{R}^{n}$ are geometrically distinct if $q_{1}(\mathbb{R}) \neq q_{2}(\mathbb{R})$. We
denote by $\mathcal{O}(\Omega, V)$ and $\widetilde{\mathcal{O}}(\Omega)$ the sets of all brake orbits and geometrically distinct brake orbits in $\bar{\Omega}$, respectively.

Remark 1.4. It is well known that via

$$
H(p, q)=\frac{1}{2}|p|^{2}+V(q)
$$

$x=(p, q)$ and $p=\dot{q}$, the elements in $\mathcal{O}(\Omega, V)$ and the solutions of (1.2)-(1.4) are one-to-one correspondent.
Definition 1.5. For $\Sigma \in \mathcal{H}_{b}^{s, c}(2 n)$, a brake orbit $(\tau, x)$ on $\Sigma$ is called symmetric if $x(\mathbb{R})=-x(\mathbb{R})$. Similarly, for a $C^{2}$ convex symmetric bounded domain $\Omega \subset$ $\mathbb{R}^{n}$, a brake orbit $(\tau, q) \in \mathcal{O}(\Omega, V)$ is called symmetric if $q(\mathbb{R})=-q(\mathbb{R})$.

Note that a brake orbit $(\tau, x) \in \mathcal{J}_{b}(\Sigma, H)$ with minimal period $\tau$ is symmetric if $x(t+\tau / 2)=-x(t)$ for $t \in \mathbb{R}$, and a brake orbit $(\tau, q) \in \mathcal{O}(\Omega, V)$ with minimal period $\tau$ is symmetric if $q(t+\tau / 2)=-q(t)$ for $t \in \mathbb{R}$.

Since 1948, many studies have been carried out for the brake orbit problem. In 1978, S. Bolotin proved in [4] the existence of brake orbits in a general setting. K. Hayashi in [12], H. Gluck and W. Ziller in [10], and V. Benci in [2] proved $\# \widetilde{\mathcal{O}}(\Omega) \geq 1$ if $V$ is $C^{1}, \bar{\Omega}=\{V \leq h\}$ is compact, and $V^{\prime}(q) \neq 0$ for all $q \in \partial \Omega$. P. Rabinowitz in [25] proved that if $H$ satisfies (1.1], $\Sigma \equiv H^{-1}(h)$ is star-shaped, and $x \cdot H^{\prime}(x) \neq 0$ for all $x \in \Sigma$, then ${ }^{\#} \tilde{\mathcal{J}}_{b}(\Sigma) \geq 1$. V. Benci and F . Giannoni gave a different proof of the existence of one brake orbit in [3]. It has been pointed out in [8] that the problem of finding brake orbits is equivalent to finding orthogonal geodesic chords on a manifold with concave boundary. R. Giambò, F. Giannoni, and P. Piccione in [9] proved the existence of an orthogonal geodesic chord on a Riemannian manifold homeomorphic to a closed disk and with concave boundary.

For multiplicity of the brake problems, A. Weinstein in [30] proved a localized result: Assume $H$ satisfies (1.1). For any $h$ sufficiently close to $H\left(z_{0}\right)$ with $z_{0}$ being a nondegenerate local minimum of $H$, there exist at least $n$ geometrically distinct brake orbits on the energy surface $H^{-1}(h)$. In [5, 10], under assumptions of Seifert in [26], it was proved that the existence of at least $n$ brake orbits, while a very strong assumption on the energy integral was used to ensure that different minimax critical levels correspond to geometrically distinct brake orbits. A. Szulkin in [27] proved that ${ }^{\#} \tilde{\mathcal{J}}_{b}\left(H^{-1}(h)\right) \geq n$ if $H$ satisfies conditions in [25] of Rabinowitz and the energy hypersurface $H^{-1}(h)$ is $\sqrt{2}$-pinched. E. van Groesen in [28] and A. Ambrosetti, V. Benci, and Y. Long in [1] also proved $\# \widetilde{\mathcal{O}}(\Omega) \geq n$ under different pinching conditions. Without a pinching condition, in [21] Y. Long, C. Zhu, and the second author of this paper proved that: For any $\Sigma \in \mathcal{H}_{b}^{s, c}(2 n)$ with $n \geq 2$, \# $\widetilde{\mathcal{J}}_{b}(\Sigma) \geq 2$. The authors of this paper in [17] proved that ${ }^{\#} \widetilde{\mathcal{J}}_{b}(\Sigma) \geq\left[\frac{n}{2}\right]+1$ for $\Sigma \in \mathcal{H}_{b}^{s, c}(2 n)$. Moreover, it was proved that if all brake orbits on $\Sigma$ are nondegenerate, then ${ }^{\#} \tilde{\mathcal{J}}_{b}(\Sigma) \geq n+\mathfrak{A}(\Sigma)$, where $2 \mathfrak{A}(\Sigma)$ is the number of geometrically distinct asymmetric brake orbits on $\Sigma$. Recently, in [34] the authors of this paper improved the results of [17] to ${ }^{\#} \widetilde{\mathcal{J}}_{b}(\Sigma) \geq\left[\frac{n+1}{2}\right]+1$ for $\Sigma \in \mathcal{H}_{b}^{s, c}(2 n)$,
$n \geq 3$. In [33] the authors of this paper proved that ${ }^{\#} \widetilde{\mathcal{J}}_{b}(\Sigma) \geq\left[\frac{n+1}{2}\right]+2$ for $\Sigma \in \mathcal{H}_{b}^{s, c}(2 n), n \geq 4$.

### 1.3 Some Consequences of Theorem 1.2 and Further Arguments

As direct consequences of Theorem 1.2 , we have the following two important corollaries:

Corollary 1.6. If $H(p, q)$ defined by (1.11) is even and convex, then the Seifert conjecture holds.

Remark 1.7. If the function $H$ in Remark 1.3 is convex and even, then $V$ is convex and even, and $\Omega$ is convex and central symmetric. Hence $\Omega$ is homeomorphic to the unit open ball in $\mathbb{R}^{n}$.

Corollary 1.8. Suppose $V(0)=0, V(q) \geq 0, V(-q)=V(q)$, and $V^{\prime \prime}(q)$ is positive definite for all $q \in \mathbb{R}^{n} \backslash\{0\}$. Then for any given $h>0$ and $\Omega \equiv\{q \in$ $\left.\mathbb{R}^{n} \mid V(q)<h\right\}$, there holds

$$
\# \widetilde{\mathcal{O}}(\Omega) \geq n
$$

It is interesting to ask the following question: Are all closed characteristics on any hypersurfaces $\Sigma \in \mathcal{H}_{b}^{s, c}(2 n)$ symmetric brake orbits after suitable time translation provided that ${ }^{\#} \widetilde{\mathcal{J}}(\Sigma)<+\infty$ ? In this direction, we have the following result:

THEOREM 1.9. For any $\Sigma \in \mathcal{H}_{b}^{s, c}(2 n)$, suppose

$$
\# \widetilde{\mathcal{J}}(\Sigma)=n
$$

Then all of the $n$ closed characteristics on $\Sigma$ are symmetric brake orbits after suitable time translation.

For $n=2$, it was proved in [13] that ${ }^{\#} \tilde{\mathcal{J}}(\Sigma)$ is either 2 or $+\infty$ for any $C^{2}$ compact convex hypersurface $\Sigma$ in $\mathbb{R}^{4}$. Hence Theorem 1.9 gives a positive answer to the above question in the case $n=2$. We also note that for the hypersurface

$$
\Sigma=\left\{\left(x_{1}, x_{2}, y_{1}, y_{2}\right) \in \mathbb{R}^{4} \left\lvert\, x_{1}^{2}+y_{1}^{2}+\frac{x_{2}^{2}+y_{2}^{2}}{4}=1\right.\right\}
$$

we have ${ }^{\#} \widetilde{\mathcal{J}}_{b}(\Sigma)=+\infty$ and ${ }^{\#} \widetilde{\mathcal{J}}_{b}^{s}(\Sigma)=2$, where we have denoted by $\widetilde{\mathcal{J}}_{b}^{s}(\Sigma)$ the set of all symmetric brake orbits on $\Sigma$. We also note that on the hypersurface $\Sigma=\left\{x \in \mathbb{R}^{2 n}| | x \mid=1\right\}$ there are some non-brake-closed characteristics.

The key ingredients in the proof of Theorem 1.2 are some ideas from our previous paper [17] and the following result, which generalizes corresponding results of our previous papers [33,34] completely, where the iteration path $\gamma^{2}$ will be defined in Definition 2.9 below.

THEOREM 1.10. For $\gamma \in \mathcal{P}_{\tau}(2 n)$, let $P=\gamma(\tau)$. If $i_{L_{0}}(\gamma) \geq 0, i_{L_{1}}(\gamma) \geq 0$, $i(\gamma) \geq n$, and $\gamma^{2}(t)=\gamma(t-\tau) \gamma(\tau)$ for all $t \in[\tau, 2 \tau]$, then

$$
\begin{equation*}
i_{L_{1}}(\gamma)+S_{P^{2}}^{+}(1)-v_{L_{0}}(\gamma) \geq 0 \tag{1.17}
\end{equation*}
$$

In this paper, we denote by $\mathbb{N}, \mathbb{Z}, \mathbb{Q}, \mathbb{R}$, and $\mathbb{C}$ the sets of positive integers, integers, rational numbers, real numbers, and complex numbers, respectively. We denote by both $\langle\cdot, \cdot\rangle$ and $\cdot$ the standard inner product in $\mathbb{R}^{n}$ or $\mathbb{R}^{2 n}$, and by $(\cdot, \cdot)$ the inner product of corresponding Hilbert space. For any $a \in \mathbb{R}$, we denote by $[a]=\max \{k \in \mathbb{Z} \mid k \leq a\}$.

We prove Theorem 1.2 and Theorem 1.9 in Section 3, and the proof of Theorem 1.10 is given in Sections 4 and 5 .

## 2 Index Theories for Symplectic Paths and the Homotopic Properties of Symplectic Matrices

In this section we make some preparations for the proof of Theorems 1.2 and 1.9 . We first briefly introduce the Maslov-type index theory of $\left(i_{L_{j}}, v_{L_{j}}\right)$ for $j=0,1$ and $\left(i_{\omega}, v_{\omega}\right)$ for $\omega \in \mathbf{U}:=\{z \in \mathbb{C}| | z \mid=1\}$.

Let $\mathcal{L}\left(\mathbb{R}^{2 n}\right)$ denote the set of $2 n \times 2 n$ real matrices and $\mathcal{L}_{S}\left(\mathbb{R}^{2 n}\right)$ its subset of symmetric ones. For any $F \in \mathcal{L}_{S}\left(\mathbb{R}^{2 n}\right)$, we denote by $m^{*}(F)$ the dimension of maximal positive definite subspace, negative definite subspace, and kernel of any $F$ for $*=+,-, 0$, respectively.

Let

$$
J_{k}=\left(\begin{array}{cc}
0 & -I_{k} \\
I_{k} & 0
\end{array}\right) \quad \text { and } \quad N_{k}=\left(\begin{array}{cc}
-I_{k} & 0 \\
0 & I_{k}
\end{array}\right)
$$

with $I_{k}$ being the identity in $\mathbb{R}^{k}$. If $k=n$ we will omit the subscript $k$ for convenience, i.e., $J_{n}=J$ and $N_{n}=N$.

The symplectic group $\operatorname{Sp}(2 k)$ for any $k \in \mathbb{N}$ is defined by

$$
\operatorname{Sp}(2 k)=\left\{M \in \mathcal{L}\left(\mathbb{R}^{2 k}\right) \mid M^{\top} J_{k} M=J_{k}\right\}
$$

where $M^{\top}$ is the transpose of matrix $M$.
For any $\tau>0$, the symplectic path in $\operatorname{Sp}(2 k)$ starting from the identity $I_{2 k}$ is defined by

$$
\mathcal{P}_{\tau}(2 k)=\left\{\gamma \in C([0, \tau], \operatorname{Sp}(2 k)) \mid \gamma(0)=I_{2 k}\right\}
$$

The Maslov-type index theory of $(i(\gamma), \nu(\gamma))$ of $\gamma$ usually plays an important role in the study of periodic solutions of Hamiltonian systems. It was introduced by C. Conley and E. Zehnder in [7] for nondegenerate symplectic path $\gamma \in \mathcal{P}_{\tau}(2 n)$ with $n \geq 2$. Y. Long and E . Zehnder in [23] extended the definition to include $\gamma \in \mathcal{P}_{\tau}(2)$. Long in [18] and C. Viterbo in [29] further extended the definition for $\gamma \in \mathcal{P}(2 n)$. In [19], Long introduced the $\omega$-index, which is an index function $\left(i_{\omega}(\gamma), v_{\omega}(\gamma)\right) \in \mathbb{Z} \times\{0,1, \ldots, 2 n\}$ for $\omega \in \mathbf{U}$ (see [20] and [18]).

For any $\omega \in \mathbf{U}$, the following hypersurface in $\operatorname{Sp}(2 n)$ is defined by

$$
\operatorname{Sp}(2 n)_{\omega}^{0}=\left\{M \in \operatorname{Sp}(2 n) \mid \operatorname{det}\left(M-\omega I_{2 n}\right)=0\right\}
$$

For any two continuous paths $\xi$ and $\eta:[0, \tau] \rightarrow \operatorname{Sp}(2 n)$ with $\xi(\tau)=\eta(0)$, their joint path is defined by

$$
\eta * \xi(t)= \begin{cases}\xi(2 t) & \text { if } 0 \leq t \leq \frac{\tau}{2}  \tag{2.1}\\ \eta(2 t-\tau) & \text { if } \frac{\tau}{2} \leq t \leq \tau\end{cases}
$$

Given any two $\left(2 m_{k} \times 2 m_{k}\right)$ matrices of square block form

$$
M_{k}=\left(\begin{array}{ll}
A_{k} & B_{k} \\
C_{k} & D_{k}
\end{array}\right)
$$

for $k=1,2$, as in [20], the $\diamond$-product (or symplectic direct product) of $M_{1}$ and $M_{2}$ is defined by the following $\left(2\left(m_{1}+m_{2}\right) \times 2\left(m_{1}+m_{2}\right)\right)$ matrix $M_{1} \diamond M_{2}$ :

$$
M_{1} \diamond M_{2}=\left(\begin{array}{cccc}
A_{1} & 0 & B_{1} & 0 \\
0 & A_{2} & 0 & B_{2} \\
C_{1} & 0 & D_{1} & 0 \\
0 & C_{2} & 0 & D_{2}
\end{array}\right) .
$$

We denote by $M^{\diamond k}$ the $k$-times self $\diamond$-product of $M$ for any $k \in \mathbb{N}$.
It is easy to see that

$$
\begin{align*}
& N_{m_{1}+m_{2}}\left(M_{1} \diamond M_{2}\right)^{-1} N_{m_{1}+m_{2}}\left(M_{1} \diamond M_{2}\right)=  \tag{2.2}\\
& \quad\left(N_{m_{1}} M_{1}^{-1} N_{m_{1}} M_{1}\right) \diamond\left(N_{m_{2}} M_{2}^{-1} N_{m_{2}} M_{2}\right) .
\end{align*}
$$

A special path $\xi_{n}$ is defined by

$$
\xi_{n}(t)=\left(\begin{array}{cc}
2-\frac{t}{\tau} & 0 \\
0 & \left(2-\frac{t}{\tau}\right)^{-1}
\end{array}\right)^{\diamond n}, \quad \forall t \in[0, \tau] .
$$

Definition 2.1. For any $\omega \in \mathbf{U}$ and $M \in \operatorname{Sp}(2 n)$, define

$$
\begin{equation*}
v_{\omega}(M)=\operatorname{dim}_{\mathbb{C}} \operatorname{ker}\left(M-\omega I_{2 n}\right) . \tag{2.3}
\end{equation*}
$$

For any $\gamma \in \mathcal{P}_{\tau}(2 n)$, define

$$
\begin{equation*}
v_{\omega}(\gamma)=v_{\omega}(\gamma(\tau)) . \tag{2.4}
\end{equation*}
$$

If $\gamma(\tau) \notin \operatorname{Sp}(2 n)_{\omega}^{0}$, we define

$$
\begin{equation*}
i_{\omega}(\gamma)=\left[\operatorname{Sp}(2 n)_{\omega}^{0}: \gamma * \xi_{n}\right], \tag{2.5}
\end{equation*}
$$

where the right-hand side of (2.5) is the usual homotopy intersection number and the orientation of $\gamma * \xi_{n}$ is its positive time direction under homotopy with fixed endpoints. If $\omega=1$, we will simply write $i(\gamma)$ instead of $i_{1}(\gamma)$. If $\gamma(\tau) \in$ $\operatorname{Sp}(2 n)_{\omega}^{0}$, we let $\mathcal{F}(\gamma)$ be the set of all open neighborhoods of $\gamma$ in $\mathcal{P}_{\tau}(2 n)$, and define

$$
\begin{equation*}
i_{\omega}(\gamma)=\sup _{U \in \mathcal{F}(\gamma)} \inf \left\{i_{\omega}(\beta) \mid \beta(\tau) \in U \text { and } \beta(\tau) \notin \operatorname{Sp}(2 n)_{\omega}^{0}\right\} \tag{2.6}
\end{equation*}
$$

For any $M \in \operatorname{Sp}(2 n)$ we define

$$
\begin{align*}
& \Omega(M)=\{P \in \operatorname{Sp}(2 n) \mid \sigma(P) \cap \mathbf{U}=\sigma(M) \cap \mathbf{U} \\
& \text { and } \left.\nu_{\lambda}(P)=\nu_{\lambda}(M) \forall \lambda \in \sigma(M) \cap \mathbf{U}\right\}, \tag{2.7}
\end{align*}
$$

where we denote by $\sigma(P)$ the spectrum of $P$.
We denote by $\Omega^{0}(M)$ the path-connected component of $\Omega(M)$ containing $M$, and call it the homotopy component of $M$ in $\operatorname{Sp}(2 n)$.

DEFINITION 2.2. For any $M_{1}, M_{2} \in \operatorname{Sp}(2 n)$, we call $M_{1} \approx M_{2}$ if $M_{1} \in \Omega^{0}\left(M_{2}\right)$.
Remark 2.3. It is easy to check that $\approx$ is an equivalence relation. If $M_{1} \approx M_{2}$, we have $M_{1}^{k} \approx M_{2}^{k}$ for any $k \in \mathbb{N}$ and $M_{1} \diamond M_{3} \approx M_{2} \diamond M_{4}$ for $M_{3} \approx M_{4}$. Also we have $M_{1} \diamond M_{2} \approx M_{2} \diamond M_{1}$ and $P M P^{-1} \approx M$ for any $P, M \in \operatorname{Sp}(2 n)$. By theorem 7.8 of [19], $M_{1} \diamond M_{2} \approx M_{1} \diamond M_{3}$ if and only if $M_{2} \approx M_{3}$.

Lemma 2.4. Assume $M_{1} \in \operatorname{Sp}\left(2\left(k_{1}+k_{2}\right)\right)$ and $M_{2} \in \operatorname{Sp}\left(2 k_{3}\right)$ have the following block form:

$$
M_{1}=\left(\begin{array}{llll}
A_{1} & A_{2} & B_{1} & B_{2} \\
A_{3} & A_{4} & B_{3} & B_{4} \\
C_{1} & C_{2} & D_{1} & D_{2} \\
C_{3} & C_{4} & D_{3} & D_{4}
\end{array}\right) \quad \text { and } \quad M_{2}=\left(\begin{array}{ll}
A_{5} & B_{5} \\
C_{5} & D_{5}
\end{array}\right)
$$

with $A_{1}, B_{1}, C_{1}, D_{1} \in \mathcal{L}\left(\mathbb{R}^{k_{1}}\right), A_{4}, B_{4}, C_{4}, D_{4} \in \mathcal{L}\left(\mathbb{R}^{k_{1}}\right)$, and $A_{5}, D_{5} \in \mathcal{L}\left(\mathbb{R}^{k_{3}}\right)$. Let

$$
M_{3}=\left(\begin{array}{cccccc}
A_{1} & 0 & A_{2} & B_{1} & 0 & B_{2} \\
0 & A_{5} & 0 & 0 & B_{5} & 0 \\
A_{3} & 0 & A_{4} & B_{3} & 0 & B_{4} \\
C_{1} & 0 & C_{2} & D_{1} & 0 & D_{2} \\
0 & C_{5} & 0 & 0 & D_{5} & 0 \\
C_{3} & 0 & C_{4} & D_{3} & 0 & D_{4}
\end{array}\right) .
$$

Then

$$
\begin{equation*}
M_{3} \approx M_{1} \diamond M_{2} . \tag{2.8}
\end{equation*}
$$

Proof. Let

$$
P=\operatorname{diag}\left(\left(\begin{array}{ccc}
I_{k_{1}} & 0 & 0 \\
0 & 0 & I_{k_{2}} \\
0 & I_{k_{3}} & 0
\end{array}\right),\left(\begin{array}{ccc}
I_{k_{1}} & 0 & 0 \\
0 & 0 & I_{k_{2}} \\
0 & I_{k_{3}} & 0
\end{array}\right)\right) .
$$

It is easy to verify that $P \in \operatorname{Sp}\left(2\left(k_{1}+k_{2}+k_{3}\right)\right)$ and $M_{3}=P\left(M_{1} \diamond M_{2}\right) P^{-1}$. Then (2.8) holds from Remark 2.3 and the proof of Lemma 2.4 is completed.

The following symplectic matrices were introduced as basic normal forms in [20]:

$$
\begin{array}{rlrl}
D(\lambda) & =\left(\begin{array}{ll}
\lambda & 0 \\
0 & \lambda^{-1}
\end{array}\right), & & \lambda= \pm 2, \\
N_{1}(\lambda, b) & =\left(\begin{array}{ll}
\lambda & b \\
0 & \lambda
\end{array}\right), & \lambda= \pm 1, b= \pm 1,0, \\
R(\theta) & =\left(\begin{array}{cc}
\cos (\theta) & -\sin (\theta) \\
\sin (\theta) & \cos (\theta)
\end{array}\right), & & \theta \in(0, \pi) \cup(\pi, 2 \pi), \\
N_{2}(\omega, b) & =\left(\begin{array}{cc}
R(\theta) & b \\
0 & R(\theta)
\end{array}\right), & & \theta \in(0, \pi) \cup(\pi, 2 \pi),
\end{array}
$$

where

$$
b=\left(\begin{array}{ll}
b_{1} & b_{2} \\
b_{3} & b_{4}
\end{array}\right)
$$

with $b_{i} \in \mathbb{R}$ and $b_{2} \neq b_{3}$.
For any $M \in \operatorname{Sp}(2 n)$ and $\omega \in \mathbf{U}$, the splitting number of $M$ at $\omega$, defined by

$$
S_{M}^{ \pm}(\omega)=\lim _{\epsilon \rightarrow 0^{+}} i_{\omega \exp ( \pm \sqrt{-1} \epsilon)}(\gamma)-i_{\omega}(\gamma)
$$

for any path $\gamma \in \mathcal{P}_{\tau}(2 n)$ satisfying $\gamma(\tau)=M$, possesses the following properties:
Lemma 2.5 ([19], [20, lemma 9.1.5 and list 9.1.12]). Splitting numbers $S_{M}^{ \pm}(\omega)$ are well-defined; i.e., they are independent of the choice of the path $\gamma \in \mathcal{P}_{\tau}(2 n)$ satisfying $\gamma(\tau)=M$. For $\omega \in \mathbf{U}$ and $M \in \operatorname{Sp}(2 n), S_{Q}^{ \pm}(\omega)=S_{M}^{ \pm}(\omega)$ if $Q \approx M$. Moreover, we have the following:
(1) $\left(S_{M}^{+}( \pm 1), S_{M}^{-}( \pm 1)\right)=(1,1)$ for $M= \pm N_{1}(1, b)$ with $b=1$ or 0 .
(2) $\left(S_{M}^{+}( \pm 1), S_{M}^{-}( \pm 1)\right)=(0,0)$ for $M= \pm N_{1}(1, b)$ with $b=-1$.
(3) $\left(S_{M}^{+}\left(e^{\sqrt{-1} \theta}\right), S_{M}^{-}\left(e^{\sqrt{-1} \theta}\right)\right)=(0,1)$ for $M=R(\theta)$ with $\theta \in(0, \pi) \cup$ $(\pi, 2 \pi)$.
(4) $\left(S_{M}^{+}(\omega), S_{M}^{-}(\omega)\right)=(0,0)$ for $\omega \in \mathbf{U} \backslash \mathbb{R}$ and $M=N_{2}(\omega, b)$ is trivial, i.e., for sufficiently small $\alpha>0, M R((t-1) \alpha)^{\diamond n}$ possesses no eigenvalues on $\mathbf{U}$ for $t \in[0,1)$.
(5) $\left(S_{M}^{+}(\omega), S_{M}^{-}(\omega)=(1,1)\right.$ for $\omega \in \mathbf{U} \backslash \mathbb{R}$ and $M=N_{2}(\omega, b)$ is nontrivial.
(6) $\left(S_{M}^{+}(\omega), S_{M}^{-}(\omega)=(0,0)\right.$ for any $\omega \in \mathbf{U}$ and $M \in \operatorname{Sp}(2 n)$ with $\sigma(M) \cap$ $\mathbf{U}=\varnothing$.
(7) $S_{M_{1} \diamond M_{2}}^{ \pm}(\omega)=S_{M_{1}}^{ \pm}(\omega)+S_{M_{2}}^{ \pm}(\omega)$ for any $M_{j} \in \operatorname{Sp}\left(2 n_{j}\right)$ with $j=1,2$ and $\omega \in \mathbf{U}$.

We denote by

$$
\begin{equation*}
F=\mathbb{R}^{2 n} \oplus \mathbb{R}^{2 n} \tag{2.9}
\end{equation*}
$$

equipped with the standard inner product $(\cdot, \cdot)$ and define the symplectic structure of $F$ by

$$
\{v, w\}=(\mathcal{J} v, w) \quad \forall v, w \in F \text { where } \mathcal{J}=(-J) \oplus J=\left(\begin{array}{rr}
-J & 0  \tag{2.10}\\
0 & J
\end{array}\right) .
$$

We denote by $\operatorname{Lag}(F)$ the set of Lagrangian subspaces of $F$ and equip it with the topology as a subspace of the Grassmannian of all $2 n$-dimensional subspaces of $F$.

It is easy to check that, for any $M \in \operatorname{Sp}(2 n)$ its graph

$$
\operatorname{Gr}(M) \equiv\left\{\left.\binom{x}{M x} \right\rvert\, x \in \mathbb{R}^{2 n}\right\}
$$

is a Lagrangian subspace of $F$.
Let

$$
\begin{align*}
& V_{1}=L_{0} \times L_{0}=\{0\} \times \mathbb{R}^{n} \times\{0\} \times \mathbb{R}^{n} \subset \mathbb{R}^{4 n},  \tag{2.11}\\
& V_{2}=L_{1} \times L_{1}=\mathbb{R}^{n} \times\{0\} \times \mathbb{R}^{n} \times\{0\} \subset \mathbb{R}^{4 n} . \tag{2.12}
\end{align*}
$$

By proposition 6.1 of [22] and lemma 2.8 and definition 2.5 of [21], we give the following:

Definition 2.6. For any continuous path $\gamma \in \mathcal{P}_{\tau}(2 n)$, we define the following Maslov-type indices:

$$
\begin{align*}
i_{L_{0}}(\gamma) & =\mu_{F}^{\mathrm{CLM}}\left(V_{1}, \operatorname{Gr}(\gamma),[0, \tau]\right)-n,  \tag{2.13}\\
i_{L_{1}}(\gamma) & =\mu_{F}^{\mathrm{CLM}}\left(V_{2}, \operatorname{Gr}(\gamma),[0, \tau]\right)-n,  \tag{2.14}\\
v_{L_{j}}(\gamma) & =\operatorname{dim}\left(\gamma(\tau) L_{j} \cap L_{j}\right), \quad j=0,1, \tag{2.15}
\end{align*}
$$

where we denote by $i_{F}^{\mathrm{CLM}}(V, W,[a, b])$ the Maslov index for Lagrangian subspace path pair ( $V, W$ ) in $F$ on $[a, b]$ defined by Cappell, Lee, and Miller in [6]. For any $M \in \operatorname{Sp}(2 n)$ and $j=0$, 1 , we also denote by $\nu_{L_{j}}(M)=\operatorname{dim}\left(M L_{j} \cap L_{j}\right)$.

The index $i_{L}(\gamma)$ for any Lagrangian subspace $L \subset \mathbb{R}^{2 n}$ and symplectic path $\gamma \in \mathcal{P}_{\tau}(2 n)$ was defined by the first author of this paper in [15] in a different way (see also [14,21]).
Definition 2.7. Let $\gamma_{0}, \gamma_{1} \in \mathcal{P}_{\tau}(2 n)$ and $j=0,1$. The paths are called $L_{j}-$ homotopic, denoted by $\gamma_{0} \sim_{L_{j}} \gamma_{1}$, if there is a map $\delta:[0,1] \rightarrow \mathcal{P}(2 n)$ such that $\delta(0)=\gamma_{0}$ and $\delta(1)=\gamma_{1}$, and $\nu_{L_{j}}(\delta(s))$ is constant for $s \in[0,1]$.
Lemma 2.8 ([15]).
(1) If $\gamma_{0} \sim_{L_{j}} \gamma_{1}$, then

$$
i_{L_{j}}\left(\gamma_{0}\right)=i_{L_{j}}\left(\gamma_{1}\right), \quad \nu_{L_{j}}\left(\gamma_{0}\right)=v_{L_{j}}\left(\gamma_{1}\right) .
$$

(2) If $\gamma=\gamma_{1} \diamond \gamma_{2} \in \mathcal{P}(2 n)$, and correspondingly $L_{j}=L_{j}^{\prime} \oplus L_{j}^{\prime \prime}$, then

$$
i_{L_{j}}(\gamma)=i_{L_{j}^{\prime}}\left(\gamma_{1}\right)+i_{L_{j}^{\prime \prime}}\left(\gamma_{2}\right), \quad v_{L_{j}}(\gamma)=v_{L_{j}^{\prime}}\left(\gamma_{1}\right)+v_{L_{j}^{\prime \prime}}\left(\gamma_{2}\right) .
$$

(3) If $\gamma \in \mathcal{P}(2 n)$ is the fundamental solution of

$$
\dot{x}(t)=J B(t) x(t)
$$

with symmetric matrix function

$$
B(t)=\left(\begin{array}{ll}
b_{11}(t) & b_{12}(t) \\
b_{21}(t) & b_{22}(t)
\end{array}\right)
$$

satisfying $b_{22}(t)>0$ for any $t \in R$, then

$$
i_{L_{0}}(\gamma)=\sum_{0<s<1} v_{L_{0}}\left(\gamma_{s}\right), \quad \gamma_{s}(t)=\gamma(s t)
$$

(4) If $b_{11}(t)>0$ for any $t \in \mathbb{R}$, then

$$
i_{L_{1}}(\gamma)=\sum_{0<s<1} v_{L_{1}}\left(\gamma_{s}\right), \quad \gamma_{s}(t)=\gamma(s t)
$$

DEFINITION 2.9. For any $\gamma \in \mathcal{P}_{\tau}$ and $k \in \mathbb{N}$, in this paper the $k$-time iteration $\gamma^{k}$ of $\gamma \in \mathcal{P}_{\tau}(2 n)$ in the brake orbit boundary sense is defined by $\left.\widetilde{\gamma}\right|_{[0, k \tau]}$, where

$$
\tilde{\gamma}(t)=\left\{\begin{array}{l}
\gamma(t-2 j \tau)\left(N \gamma(\tau)^{-1} N \gamma(\tau)\right)^{j}, t \in[2 j \tau,(2 j+1) \tau], j=0,1, \ldots \\
N \gamma(2 j \tau+2 \tau-t) N\left(N \gamma(\tau)^{-1} N \gamma(\tau)\right)^{j+1} \\
\quad t \in[(2 j+1) \tau,(2 j+2) \tau], j=0,1, \ldots
\end{array}\right.
$$

## 3 Proofs of Theorems 1.2 and 1.9

In this section we prove Theorems 1.2 and 1.9
For $\Sigma \in \mathcal{H}_{b}^{s, c}(2 n)$, let $j_{\Sigma}: \Sigma \rightarrow[0,+\infty)$ be the gauge function of $\Sigma$ defined by

$$
j_{\Sigma}(0)=0 \quad \text { and } \quad j_{\Sigma}(x)=\inf \left\{\lambda>0 \left\lvert\, \frac{x}{\lambda} \in C\right.\right\} \quad \forall x \in \mathbb{R}^{2 n} \backslash\{0\}
$$

where $C$ is the domain enclosed by $\Sigma$.
Define

$$
\begin{equation*}
H_{\alpha}(x)=\left(j_{\Sigma}(x)\right)^{\alpha}, \quad \alpha>1, \quad H_{\Sigma}(x)=H_{2}(x) \quad \forall x \in \mathbb{R}^{2 n} \tag{3.1}
\end{equation*}
$$

Then $H_{\Sigma} \in C^{2}\left(\mathbb{R}^{2 n} \backslash\{0\}, \mathbb{R}\right) \cap C^{1,1}\left(\mathbb{R}^{2 n}, \mathbb{R}\right)$.
We consider the following fixed energy problem:

$$
\begin{align*}
\dot{x}(t) & =J H_{\Sigma}^{\prime}(x(t))  \tag{3.2}\\
H_{\Sigma}(x(t)) & =1  \tag{3.3}\\
x(-t) & =N x(t)  \tag{3.4}\\
x(\tau+t) & =x(t) \quad \forall t \in \mathbb{R} . \tag{3.5}
\end{align*}
$$

Denote by $\mathcal{J}_{b}(\Sigma, 2)\left(\mathcal{J}_{b}(\Sigma, \alpha)\right.$ for $\alpha=2$ in (3.1) the set of all solutions $(\tau, x)$ of problem $(3.2)-(3.5)$ and by $\widetilde{\mathcal{J}}_{b}(\Sigma, 2)$ the set of all geometrically distinct solutions of (3.2)-(3.5). By remark 1.2 of [17] or the discussion in [21],
elements in $\mathcal{J}_{b}(\Sigma)$ and $\mathcal{J}_{b}(\Sigma, 2)$ are in one-to-one correspondence. So we have $\# \widetilde{\mathcal{J}}_{b}(\Sigma)={ }^{\#} \tilde{\mathcal{J}}_{b}(\Sigma, 2)$.

For the readers' convenience, in the following we list some known results that will be used in the proof of Theorem 1.2 .

In the following we write $\left(i_{L_{0}}(\gamma, k), \nu_{L_{0}}(\gamma, k)\right)=\left(i_{L_{0}}\left(\gamma^{k}\right), \nu_{L_{0}}\left(\gamma^{k}\right)\right)$ for any symplectic path $\gamma \in \mathcal{P}_{\tau}(2 n)$ and $k \in \mathbb{N}$, where $\gamma^{k}$ is defined by Definition 2.9 .
Lemma 3.1 (Theorem 1.5 of [17] and Theorem 4.3 of [22]). Let $\gamma_{j} \in \mathcal{P}_{\tau_{j}}(2 n)$ for $j=1, \ldots, q$. Let $M_{j}=\widehat{\gamma_{j}^{2}}\left(2 \tau_{j}\right)=N \gamma_{j}\left(\tau_{j}\right)^{-1} N \gamma_{j}\left(\tau_{j}\right)$ for $j=1, \ldots, q$. Suppose

$$
\hat{i}_{L_{0}}\left(\gamma_{j}\right)>0, \quad j=1, \ldots, q .
$$

Then there exist infinitely many $\left(R, m_{1}, m_{2}, \ldots, m_{q}\right) \in \mathbb{N}^{q+1}$ such that
(i) $\nu_{L_{0}}\left(\gamma_{j}, 2 m_{j} \pm 1\right)=\nu_{L_{0}}\left(\gamma_{j}\right)$,
(ii) $i_{L_{0}}\left(\gamma_{j}, 2 m_{j}-1\right)+v_{L_{0}}\left(\gamma_{j}, 2 m_{j}-1\right)=R-\left(i_{L_{1}}\left(\gamma_{j}\right)+n+S_{M_{j}}^{+}(1)-\right.$ $\left.\nu_{L_{0}}\left(\gamma_{j}\right)\right)$,
(iii) $i_{L_{0}}\left(\gamma_{j}, 2 m_{j}+1\right)=R+i_{L_{0}}\left(\gamma_{j}\right)$,
(iv) $\nu\left(\gamma_{j}^{2}, 2 m_{j} \pm 1\right)=\nu\left(\gamma_{j}^{2}\right)$,
(v) $i\left(\gamma_{j}^{2}, 2 m_{j}-1\right)+v\left(\gamma_{j}^{2}, 2 m_{j}-1\right)=2 R-\left(i\left(\gamma_{j}^{2}\right)+2 S_{M_{j}}^{+}(1)-v\left(\gamma_{j}^{2}\right)\right)$,
(vi) $i\left(\gamma_{j}^{2}, 2 m_{j}+1\right)=2 R+i\left(\gamma_{j}^{2}\right)$,
where we have set $i\left(\gamma_{j}^{2}, n_{j}\right)=i\left(\gamma_{j}^{2 n_{j}}\right), \nu\left(\gamma_{j}^{2}, n_{j}\right)=v\left(\gamma_{j}^{2 n_{j}}\right)$ for $n_{j} \in \mathbb{N}$.
For any $(\tau, x) \in \mathcal{J}_{b}(\Sigma, 2)$, there is a corresponding path $\gamma_{x} \in \mathcal{P}_{\tau}(2 n)$. For $m \in$ $\mathbb{N}$, we denote by $i_{L_{j}}(x, m)=i_{L_{j}}\left(\gamma_{x}^{m}\right)$ and $\nu_{L_{j}}(x, m)=v_{L_{j}}\left(\gamma_{x}^{m}\right)$ for $j=0,1$. Also we denote $i(x, m)=i\left(\gamma_{x}^{2 m}\right)$ and $\nu(x, m)=\nu\left(\gamma_{x}^{2 m}\right)$. We remind the reader that the symplectic path $\gamma_{x}^{m}$ is defined in the interval $\left[0, \frac{m \tau}{2}\right]$, and the symplectic path $\gamma_{x}^{2 m}$ is defined in the interval $[0, m \tau]$. If $m=1$, we denote $i(x)=i(x, 1)$ and $v(x)=v(x, 1)$. By lemma 6.3 of [17] we have the following:
Lemma 3.2. Suppose ${ }^{\#} \tilde{\mathcal{J}}_{b}(\Sigma)<+\infty$. Then there exist an integer $K \geq 0$ and an injective map $\phi: \mathbb{N}+K \mapsto \mathcal{J}_{b}(\Sigma, 2) \times \mathbb{N}$ such that
(i) For any $k \in \mathbb{N}+K,[(\tau, x)] \in \mathcal{J}_{b}(\Sigma, 2)$, and $m \in \mathbb{N}$ satisfying $\phi(k)=$ ( $[(\tau, x)], m$ ), there holds

$$
i_{L_{0}}(x, m) \leq k-1 \leq i_{L_{0}}(x, m)+v_{L_{0}}(x, m)-1,
$$

where $x$ has minimal period $\tau$.
(ii) For any $k_{j} \in \mathbb{N}+K, k_{1}<k_{2}$, and $\left(\tau_{j}, x_{j}\right) \in \mathcal{J}_{b}(\Sigma, 2)$ satisfying $\phi\left(k_{j}\right)=$ $\left(\left[\left(\tau_{j}, x_{j}\right)\right], m_{j}\right)$ with $j=1,2$ and $\left[\left(\tau_{1}, x_{1}\right)\right]=\left[\left(\tau_{2}, x_{2}\right)\right]$, there holds

$$
m_{1}<m_{2} .
$$

Lemma 3.3 (Lemma 7.2 of [17]). Let $\gamma \in \mathcal{P}_{\tau}(2 n)$ be extended to $[0,+\infty)$ by $\gamma(\tau+t)=\gamma(t) \gamma(\tau)$ for all $t>0$. Suppose $\gamma(\tau)=M=P^{-1}\left(I_{2} \diamond \tilde{M}\right) P$ with
$\tilde{M} \in \operatorname{Sp}(2 n-2)$ and $i(\gamma) \geq n$. Then we have

$$
i(\gamma, 2)+2 S_{M^{2}}^{+}(1)-v(\gamma, 2) \geq n+2
$$

Lemma 3.4 (Lemma 7.3 of [17]). For any $(\tau, x) \in \mathcal{J}_{b}(\Sigma, 2)$ and $m \in \mathbb{N}$, there hold

$$
\begin{aligned}
i_{L_{0}}(x, m+1)-i_{L_{0}}(x, m) & \geq 1 \\
i_{L_{0}}(x, m+1)+v_{L_{0}}(x, m+1)-1 & \geq i_{L_{0}}(x, m+1) \\
& >i_{L_{0}}(x, m)+v_{L_{0}}(x, m)-1
\end{aligned}
$$

Proof of Theorem 1.2, It is suffices to consider the case ${ }^{\#} \widetilde{\mathcal{J}}_{b}(\Sigma)<+\infty$. Since $-\Sigma=\Sigma$, for $(\tau, x) \in \mathcal{J}_{b}(\Sigma, 2)$ we have

$$
\begin{equation*}
H_{\Sigma}(x)=H_{\Sigma}(-x), \quad H_{\Sigma}^{\prime}(x)=-H_{\Sigma}^{\prime}(-x), \quad H_{\Sigma}^{\prime \prime}(x)=H_{\Sigma}^{\prime \prime}(-x) \tag{3.6}
\end{equation*}
$$

It follows that $(\tau,-x) \in \mathcal{J}_{b}(\Sigma, 2)$ and, in view of the definition of $\gamma_{x}$, we obtain that

$$
\gamma_{x}=\gamma_{-x}
$$

Hence

$$
\begin{align*}
& \left(i_{L_{0}}(x, m), v_{L_{0}}(x, m)\right)=\left(i_{L_{0}}(-x, m), v_{L_{0}}(-x, m)\right) \\
& \left(i_{L_{1}}(x, m), v_{L_{1}}(x, m)\right)=\left(i_{L_{1}}(-x, m), v_{L_{1}}(-x, m)\right)
\end{align*} \quad \forall m \in \mathbb{N}
$$

We can write

$$
\begin{align*}
\widetilde{\mathcal{J}}_{b}(\Sigma, 2)= & \left\{\left[\left(\tau_{j}, x_{j}\right)\right] \mid j=1, \ldots, p\right\}  \tag{3.8}\\
& \cup\left\{\left[\left(\tau_{k}, x_{k}\right)\right],\left[\left(\tau_{k},-x_{k}\right)\right] \mid k=p+1, \ldots, p+q\right\}
\end{align*}
$$

with $x_{j}(\mathbb{R})=-x_{j}(\mathbb{R})$ for $j=1, \ldots, p$ and $x_{k}(\mathbb{R}) \neq-x_{k}(\mathbb{R})$ for $k=p+$ $1, \ldots, p+q$. Here we recall that $\left(\tau_{j}, x_{j}\right)$ has minimal period $\tau_{j}$ for $j=1, \ldots$, $p+q$ and $x_{j}\left(\frac{\tau_{j}}{2}+t\right)=-x_{j}(t), t \in \mathbb{R}$, for $j=1, \ldots, p$.

In view of Lemma 3.2 there exists an integer $K \geq 0$ and an injective map $\phi: \mathbb{N}+K \rightarrow \mathcal{J}_{b}(\Sigma, 2) \times \mathbb{N}$. By (3.7), $\left(\tau_{k}, x_{k}\right)$ and $\left(\tau_{k},-x_{k}\right)$ have the same $\left(i_{L_{0}}, v_{L_{0}}\right)$-indices. So by Lemma 3.2, without loss of generality, we can further require that

$$
\begin{equation*}
\operatorname{Im}(\phi) \subseteq\left\{\left[\left(\tau_{k}, x_{k}\right)\right] \mid k=1, \ldots, p+q\right\} \times \mathbb{N} \tag{3.9}
\end{equation*}
$$

By the strict convexity of $H_{\Sigma}$ and (6.19) of [17]), we have

$$
\hat{i}_{L_{0}}\left(x_{k}\right)>0, \quad k=1, \ldots, p+q
$$

Applying Lemma 3.1 to symplectic paths

$$
\gamma_{1}, \ldots, \gamma_{p+q}, \gamma_{p+q+1}, \ldots, \gamma_{p+2 q}
$$

associated with $\left(\tau_{1}, x_{1}\right), \ldots,\left(\tau_{p+q}, x_{p+q}\right),\left(2 \tau_{p+1}, x_{p+1}^{2}\right), \ldots,\left(2 \tau_{p+q}, x_{p+q}^{2}\right)$, respectively, there exists a vector $\left(R, m_{1}, \ldots, m_{p+2 q}\right) \in \mathbb{N}^{p+2 q+1}$ such that $R>$ $K+n$ and

$$
\begin{equation*}
i_{L_{0}}\left(x_{k}, 2 m_{k}+1\right)=R+i_{L_{0}}\left(x_{k}\right) \tag{3.10}
\end{equation*}
$$

$$
\begin{align*}
& i_{L_{0}}\left(x_{k}, 2 m_{k}-1\right)+v_{L_{0}}\left(x_{k}, 2 m_{k}-1\right)=  \tag{3.11}\\
& \quad R-\left(i_{L_{1}}\left(x_{k}\right)+n+S_{M_{k}}^{+}(1)-v_{L_{0}}\left(x_{k}\right)\right)
\end{align*}
$$

for $k=1, \ldots, p+q, M_{k}=\gamma_{k}^{2}\left(\tau_{k}\right)$, and

$$
\begin{align*}
i_{L_{0}}\left(x_{k}, 4 m_{k}-2\right)+v_{L_{0}} & \left(x_{k}, 4 m_{k}-2\right)=  \tag{3.13}\\
& R-\left(i_{L_{1}}\left(x_{k}, 2\right)+n+S_{M_{k}}^{+}(1)-v_{L_{0}}\left(x_{k}, 2\right)\right)
\end{align*}
$$

for $k=p+q+1, \ldots, p+2 q$ and $M_{k}=\gamma_{k}^{4}\left(2 \tau_{k}\right)=\gamma_{k}^{2}\left(\tau_{k}\right)^{2}$.
By Lemma 3.1, we also have

$$
\begin{equation*}
i\left(x_{k}, 2 m_{k}+1\right)=2 R+i\left(x_{k}\right), \tag{3.14}
\end{equation*}
$$

for $k=1, \ldots, p+q, M_{k}=\gamma_{k}^{2}\left(\tau_{k}\right)$, and

$$
\begin{align*}
i\left(x_{k}, 4 m_{k}+2\right) & =2 R+i\left(x_{k}, 2\right),  \tag{3.16}\\
i\left(x_{k}, 4 m_{k}-2\right)+v\left(x_{k}, 4 m_{k}-2\right) & =  \tag{3.17}\\
2 R & -\left(i\left(x_{k}, 2\right)+2 S_{M_{k}}^{+}(1)-v\left(x_{k}, 2\right)\right),
\end{align*}
$$

for $k=p+q+1, \ldots, p+2 q$ and $M_{k}=\gamma_{k}^{4}\left(2 \tau_{k}\right)=\gamma_{k}^{2}\left(\tau_{k}\right)^{2}$.
From (3.9), we can set

$$
\phi(R-(s-1))=\left(\left[\left(\tau_{k(s)}, x_{k(s)}\right)\right], m(s)\right) \quad \forall s \in S:=\{1, \ldots, n\}
$$

where $k(s) \in\{1, \ldots, p+q\}$ and $m(s) \in \mathbb{N}$.
We continue our proof to study the symmetric and asymmetric orbits separately. Let

$$
S_{1}=\{s \in S \mid k(s) \leq p\}, \quad S_{2}=S \backslash S_{1} .
$$

We shall prove that ${ }^{\#} S_{1} \leq p$ and ${ }^{\#} S_{2} \leq 2 q$. These estimates together with the definitions of $S_{1}$ and $S_{2}$ yield Theorem 1.2 .
Claim 3.5. ${ }^{\#} S_{1} \leq p$.
Proof. By the definition of $S_{1}$, we have that $\left(\left[\left(\tau_{k(s)}, x_{k(s)}\right)\right], m(s)\right)$ is symmetric when $k(s) \leq p$. We further prove that $m(s)=2 m_{k(s)}$ for $s \in S_{1}$.

In fact, by the definition of $\phi$ and Lemma 3.2, for all $s=1, \ldots, n$ we have

$$
\begin{aligned}
i_{L_{0}}\left(x_{k(s)}, m(s)\right) & \leq(R-(s-1))-1=R-s \\
& \leq i_{L_{0}}\left(x_{k(s)}, m(s)\right)+v_{L_{0}}\left(x_{k(s)}, m(s)\right)-1 .
\end{aligned}
$$

By the strict convexity of $H_{\Sigma}$ and Lemma 2.8 , we have $i_{L_{0}}\left(x_{k(s)}\right) \geq 0$, so that

$$
\begin{align*}
i_{L_{0}}\left(x_{k(s)}, m(s)\right) \leq R-s<R & \leq R+i_{L_{0}}\left(x_{k(s)}\right) \\
& =i_{L_{0}}\left(x_{k(s)}, 2 m_{k(s)}+1\right) \tag{3.18}
\end{align*}
$$

for every $s=1, \ldots, n$, where we have used (3.10) in the last equality. Note that the proofs of $(3.18)$ and $(3.18)$ do not depend on the condition $s \in S_{1}$.

It is easy to see that $\gamma_{x_{k}}$ satisfies the conditions of Theorem 1.10 with $\tau=\tau_{k} / 2$. Note that by definition $i_{L_{1}}\left(x_{k}\right)=i_{L_{1}}\left(\gamma_{x_{k}}\right)$ and $\nu_{L_{0}}\left(x_{k}\right)=\nu_{L_{0}}\left(\gamma_{x_{k}}\right)$. So by Theorem 1.10 we have

$$
\begin{equation*}
i_{L_{1}}\left(x_{k}\right)+S_{M_{k}}^{+}(1)-v_{L_{0}}\left(x_{k}\right) \geq 0 \quad \forall k=1, \ldots, p \tag{3.19}
\end{equation*}
$$

Hence by (3.18) and (3.19), if $k(s) \leq p$, it follows that

$$
\begin{align*}
& i_{L_{0}}\left(r x_{k(s)}, 2 m_{k(s)}-1\right)+v_{L_{0}}\left(x_{k(s)}, 2 m_{k(s)}-1\right)-1 \\
& \quad=R-\left(i_{L_{1}}\left(x_{k(s)}\right)+n+S_{M_{k(s)}}^{+}(1)-v_{L_{0}}\left(x_{k(s)}\right)\right)-1 \\
& \quad \leq R-\frac{1-n}{2}-1-n \\
& \quad<R-s \\
& \quad \leq i_{L_{0}}\left(x_{k(s)}, m(s)\right)+v_{L_{0}}\left(x_{k(s)}, m(s)\right)-1 \tag{3.20}
\end{align*}
$$

Thus by (3.18), (3.20), and Lemma 3.4 we obtain

$$
2 m_{k(s)}-1<m(s)<2 m_{k(s)}+1 .
$$

Hence

$$
m(s)=2 m_{k(s)} \quad \text { and } \quad \phi(R-s+1)=\left(\left[\left(\tau_{k(s)}, x_{k(s)}\right)\right], 2 m_{k(s)}\right) \quad \forall s \in S_{1} .
$$

Then the injectivity of the map $\phi$ induces an injective map

$$
\phi_{1}: S_{1} \rightarrow\{1, \ldots, p\}, \quad s \mapsto k(s) .
$$

Therefore, ${ }^{\#} S_{1} \leq p$ and Claim 3.5 is proved.
Claim 3.6. ${ }^{\#} S_{2} \leq 2 q$.
Proof. By the formulas (3.14]-(3.17), and (59) of [16] (also [20, claim 4, p. 352]), we have

$$
\begin{equation*}
m_{k}=2 m_{k+q} \quad \text { for } k=p+1, p+2, \ldots, p+q . \tag{3.21}
\end{equation*}
$$

By Theorem 1.10 there holds

$$
\begin{equation*}
i_{L_{1}}\left(x_{k}, 2\right)+S_{M_{k}}^{+}(1)-v_{L_{0}}\left(x_{k}, 2\right) \geq 0, \quad p+1 \leq k \leq p+q . \tag{3.22}
\end{equation*}
$$

By (3.13), (3.18), (3.21) and (3.22), for $p+1 \leq k(s) \leq p+q$ we have

$$
\begin{aligned}
& i_{L_{0}}\left(x_{k(s)}, 2 m_{k(s)}-2\right)+v_{L_{0}}\left(x_{k(s)}, 2 m_{k(s)}-2\right)-1 \\
& \quad=i_{L_{0}}\left(x_{k(s)}, 4 m_{k(s)+q}-2\right)+v_{L_{0}}\left(x_{k(s)}, 4 m_{k(s)+q}-2\right)-1 \\
& \quad=R-\left(i_{L_{1}}\left(x_{k(s)}, 2\right)+n+S_{M_{k(s)}}^{+}(1)-v_{L_{0}}\left(x_{k(s)}, 2\right)\right)-1 \\
& \quad=R-\left(i_{L_{1}}\left(x_{k}, 2\right)+S_{M_{k}}^{+}(1)-v_{L_{0}}\left(x_{k}, 2\right)\right)-1-n \\
& \quad \leq R-1-n<
\end{aligned}
$$

$$
\begin{align*}
& <R-s \\
& \leq i_{L_{0}}\left(x_{k(s)}, m(s)\right)+v_{L_{0}}\left(x_{k(s)}, m(s)\right)-1 \tag{3.23}
\end{align*}
$$

Thus (3.18, 3.23) and Lemma 3.4 imply

$$
2 m_{k(s)}-2<m(s)<2 m_{k(s)}+1, \quad p<k(s) \leq p+q
$$

So

$$
m(s) \in\left\{2 m_{k(s)}-1,2 m_{k(s)}\right\} \quad \text { for } p<k(s) \leq p+q
$$

In particular, this yields that for any $s_{0}$ and $s \in S_{2}$, if $k(s)=k\left(s_{0}\right)$, then

$$
m(s) \in\left\{2 m_{k(s)}-1,2 m_{k(s)}\right\}=\left\{2 m_{k\left(s_{0}\right)}-1,2 m_{k\left(s_{0}\right)}\right\}
$$

Then, in view of the injectivity of the map $\phi$ from Lemma 3.2, we have

$$
\#\left\{s \in S_{2} \mid k(s)=k\left(s_{0}\right)\right\} \leq 2
$$

This proves Claim 3.6 .
By Claim 3.5 and Claim 3.6, we obtain

$$
{ }^{\#} \tilde{\mathcal{J}}_{b}(\Sigma)={ }^{\#} \tilde{\mathcal{J}}_{b}(\Sigma, 2)=p+2 q \geq{ }^{\#} S_{1}+{ }^{\#} S_{2}=n
$$

The proof of Theorem 1.2 is completed.
PROOF OF THEOREM 1.9. We call a closed characteristic $x$ on $\Sigma$ a dual brake orbit on $\Sigma$ if $x(-t)=-N x(t)$. Then by the similar proof of lemma 3.1 of [31], a closed characteristic $x$ on $\Sigma$ can become a dual brake orbit after suitable time translation if and only if $x(\mathbb{R})=-N x(\mathbb{R})$. So by lemma 3.1 of [31] again, if a closed characteristic $x$ on $\Sigma$ can both become brake orbits and dual brake orbits after suitable translation, then $x(\mathbb{R})=N x(\mathbb{R})=-N x(\mathbb{R})$. Thus $x(\mathbb{R})=-x(\mathbb{R})$.

Since we also have $-N \Sigma=\Sigma,(-N)^{2}=I_{2 n}$, and $(-N) J=-J(-N)$, dually by the same proof of Theorem 1.2 (with the estimate (5.3) in Theorem5.3 below), there are at least $n$ geometrically distinct dual brake orbits on $\Sigma$.

If there are exactly $n$ closed characteristics on $\Sigma$, then Theorem 1.2 implies that all of them are brake orbits on $\Sigma$ after suitable time translation. By the same argument all the $n$ closed characteristics must be dual brake orbits on $\Sigma$. Then by the argument in the first paragraph of the proof of this theorem, all these $n$ closed characteristics on $\Sigma$ must be symmetric. Hence all of them are symmetric brake orbits after suitable time translation. The proof of Theorem 1.9 is completed.

## $4\left(L_{0}, L_{1}\right)$-Concavity and $\left(\varepsilon, L_{0}, L_{1}\right)$-Signature of Symplectic Matrix

Definition 4.1. For any $P \in \operatorname{Sp}(2 n)$ and $\varepsilon \in \mathbb{R}$, we define the $\left(\varepsilon, L_{0}, L_{1}\right)$ symmetrization of $P$ by

$$
M_{\varepsilon}(P)=P^{\top}\left(\begin{array}{rr}
\sin 2 \varepsilon I_{n} & -\cos 2 \varepsilon I_{n} \\
-\cos 2 \varepsilon I_{n} & -\sin 2 \varepsilon I_{n}
\end{array}\right) P+\left(\begin{array}{rr}
\sin 2 \varepsilon I_{n} & \cos 2 \varepsilon I_{n} \\
\cos 2 \varepsilon I_{n} & -\sin 2 \varepsilon I_{n}
\end{array}\right)
$$

The $\left(\varepsilon, L_{0}, L_{1}\right)$-signature of $P$ is defined by the signature of $M_{\varepsilon}(P)$. The $\left(L_{0}, L_{1}\right)$ concavity and $\left(L_{0}, L_{1}\right)^{*}$-concavity of a symplectic path $\gamma$ is defined by

$$
\begin{aligned}
& \operatorname{concav}_{\left(L_{0}, L_{1}\right)}(\gamma)=i_{L_{0}}(\gamma)-i_{L_{1}}(\gamma) \\
& \operatorname{concav}_{\left(L_{0}, L_{1}\right)}^{*}(\gamma)=\left(i_{L_{0}}(\gamma)+v_{L_{0}}(\gamma)\right)-\left(i_{L_{1}}(\gamma)+v_{L_{1}}(\gamma)\right)
\end{aligned}
$$

respectively.
In [15] it was proved that $\left(L_{0}, L_{1}\right)$-concavity is only dependent on the end matrix $\gamma(\tau)$ of $\gamma$, and in [32] it was proved that the $\left(L_{0}, L_{1}\right)$-concavity of a symplectic path $\gamma$ is half of the $\left(\varepsilon, L_{0}, L_{1}\right)$-signature of $\gamma(\tau)$; i.e., we have the following result:

THEOREM 4.2 ([]32]). For $\gamma \in \mathcal{P}_{\tau}(2 k)$ with $\tau>0$, we have

$$
\operatorname{concav}_{\left(L_{0}, L_{1}\right)}(\gamma)=\frac{1}{2} \operatorname{sgn} M_{\varepsilon}(\gamma(\tau))
$$

where $0<\varepsilon \ll 1$, and we have denoted by $\operatorname{sgn} A$ the signature of $A$ for any symmetric matrix $A$. We also have

$$
\operatorname{concav}_{\left(L_{0}, L_{1}\right)}^{*}(\gamma)=\frac{1}{2} \operatorname{sgn} M_{\varepsilon}(\gamma(\tau)), \quad 0<-\varepsilon \ll 1
$$

Remark 4.3 (Remark 2.1 of [32]). For any $2 n_{j} \times 2 n_{j}$ symplectic matrix $P_{j}$ with $j=1,2$ and $n_{j} \in \mathbb{N}$, we have

$$
\begin{aligned}
M_{\varepsilon}\left(P_{1} \diamond P_{2}\right) & =M_{\varepsilon}\left(P_{1}\right) \diamond M_{\varepsilon}\left(P_{2}\right) \\
\operatorname{sgn} M_{\varepsilon}\left(P_{1} \diamond P_{2}\right) & =\operatorname{sgn} M_{\varepsilon}\left(P_{1}\right)+\operatorname{sgn} M_{\varepsilon}\left(P_{2}\right)
\end{aligned}
$$

where $\varepsilon \in \mathbb{R}$.
In the rest of this section, we further develop some basic properties of the ( $\varepsilon, L_{0}, L_{1}$ )-signature and study the normal forms of $L_{0}$-degenerate symplectic matrices.

Lemma 4.4 (Lemma 2.3 of [34]). Let $k \in \mathbb{N}$ and let

$$
P=\left(\begin{array}{cc}
I_{k} & 0 \\
C & I_{k}
\end{array}\right)
$$

be any symplectic matrix. Then $P \approx I_{2}^{\diamond p} \diamond N_{1}(1,1)^{\diamond q} \diamond N_{1}(1,-1)^{\diamond r}$ with $p=$ $m^{0}(C), q=m^{-}(C)$, and $r=m^{+}(C)$.

DEFINITION 4.5. We call two symplectic matrices $M_{1}$ and $M_{2}\left(L_{0}, L_{1}\right)$-homotopic equivalent in $\operatorname{Sp}(2 k)$, and denote the relationship by $M_{1} \sim M_{2}$, if there are $P_{j} \in \operatorname{Sp}(2 k)$ of the form $P_{j}=\operatorname{diag}\left(Q_{j},\left(Q_{j}^{\top}\right)^{-1}\right)$, where $Q_{j}$ is a $k \times k$ invertible real matrix with $\operatorname{det}\left(Q_{j}\right)>0$ for $j=1,2$ such that

$$
M_{1}=P_{1} M_{2} P_{2}
$$

Remark 4.6. Let

$$
M_{i}=\left(\begin{array}{ll}
A_{i} & B_{i} \\
C_{i} & D_{i}
\end{array}\right) \in \operatorname{Sp}\left(2 k_{i}\right), \quad i=0,1,2
$$

$k_{1}=k_{2}$, and $M_{1} \sim M_{2}$; then $A_{1}^{\top} C_{1}$ and $B_{1}^{\top} D_{1}$ are congruent to $A_{2}^{\top} C_{2}$ and $B_{2}^{\top} D_{2}$, respectively. So $m^{*}\left(A_{1}^{\top} C_{1}\right)=m^{*}\left(A_{2}^{\top} C_{2}\right)$ and $m^{*}\left(B_{1}^{\top} D_{1}\right)=m^{*}\left(B_{2}^{\top} D_{2}\right)$ for $*= \pm, 0$. Furthermore, if $M_{0}=M_{1} \diamond M_{2}$ (here $k_{1}=k_{2}$ is not necessary), then

$$
\begin{align*}
m^{*}\left(A_{0}^{\top} C_{0}\right) & =m^{*}\left(A_{1}^{\top} C_{1}\right)+m^{*}\left(A_{2}^{\top} C_{2}\right) \\
m^{*}\left(B_{0}^{\top} D_{0}\right) & =m^{*}\left(B_{1}^{\top} D_{1}\right)+m^{*}\left(B_{2}^{\top} D_{2}\right) \tag{4.1}
\end{align*}
$$

and so $m^{*}\left(A^{\top} C\right)$ and $m^{*}\left(B^{\top} D\right)$ are $\left(L_{0}, L_{1}\right)$-homotopic invariant. The following formula will be used frequently:

$$
N_{k} M_{1}^{-1} N_{k} M_{1}=I_{2 k}+2\left(\begin{array}{ll}
B_{1}^{\top} C_{1} & B_{1}^{\top} D_{1}  \tag{4.2}\\
A_{1}^{\top} C_{1} & C_{1}^{\top} B_{1}
\end{array}\right)
$$

It is clear that $\sim$ is an equivalence relation and we have the following lemma:
LEMMA 4.7 (Lemma 2.4 of [34]). For $M_{1}, M_{2} \in \operatorname{Sp}(2 k)$, if $M_{1} \sim M_{2}$, then

$$
\begin{gathered}
\operatorname{sgn} M_{\varepsilon}\left(M_{1}\right)=\operatorname{sgn} M_{\varepsilon}\left(M_{2}\right), \quad 0 \leq|\varepsilon| \ll 1 \\
N_{k} M_{1}^{-1} N_{k} M_{1} \approx N_{k} M_{2}^{-1} N_{k} M_{2}
\end{gathered}
$$

By results in [32-34], we have the following Lemmas $4.8,4.10$, which will be used frequently in Section 4

Lemma 4.8 (Lemma 2.5 of [34]). Assume

$$
P=\left(\begin{array}{ll}
A & B \\
C & D
\end{array}\right) \in \operatorname{Sp}(2 k)
$$

where $A, B, C$, and $D$ are all $k \times k$ matrices.
(i) Let $q=\max \left\{m^{+}\left(A^{\top} C\right), m^{+}\left(B^{\top} D\right)\right\}$; we have

$$
\begin{array}{ll}
\frac{1}{2} \operatorname{sgn} M_{\varepsilon}(P) \leq k-q-v_{L_{1}}(P), & 0<-\varepsilon \ll 1 \\
\frac{1}{2} \operatorname{sgn} M_{\varepsilon}(P) \leq k-q-v_{L_{0}}(P), & 0<\varepsilon \ll 1
\end{array}
$$

(ii) If both $B$ and $C$ are invertible, then

$$
\operatorname{sgn} M_{\varepsilon}(P)=\operatorname{sgn} M_{0}(P), \quad 0 \leq|\varepsilon| \ll 1
$$

LEMMA 4.9 ([]32]). For $\gamma \in \mathcal{P}_{\tau}(2), b>0$, and $\varepsilon>0$ small enough we have

$$
\operatorname{sgn} M_{ \pm \varepsilon}(R(\theta))=0 \quad \text { for } \theta \in \mathbb{R}
$$

$$
\operatorname{sgn} M_{ \pm \varepsilon}(P)=0 \quad \text { if } P=\left(\begin{array}{cc}
a & 0 \\
0 & \frac{1}{a}
\end{array}\right) \text { with } a \in \mathbb{R} \backslash\{0\}
$$

$$
\begin{array}{ll}
\operatorname{sgn} M_{\varepsilon}(P)=0 & \text { if } P= \pm\left(\begin{array}{ll}
1 & b \\
0 & 1
\end{array}\right) \text { or } \pm\left(\begin{array}{rr}
1 & 0 \\
-b & 1
\end{array}\right) \\
\operatorname{sgn} M_{\varepsilon}(P)=2 & \text { if } P= \pm\left(\begin{array}{rr}
1 & -b \\
0 & 1
\end{array}\right) \\
\operatorname{sgn} M_{\varepsilon}(P)=-2 & \text { if } P= \pm\left(\begin{array}{ll}
1 & 0 \\
b & 1
\end{array}\right)
\end{array}
$$

LEMMA 4.10 (Lemma 2.9 of [33]). Let $2 k \times 2 k$ symmetric real matrix $E$ have the block form

$$
E=\left(\begin{array}{cc}
0 & E_{1} \\
E_{1}^{\top} & E_{2}
\end{array}\right)
$$

Then

$$
\begin{equation*}
m^{ \pm}(E) \geq \operatorname{rank} E_{1} \tag{4.3}
\end{equation*}
$$

Lemma 4.13 and Lemma 4.14 are key technical results of this paper. The next lemma is used in the proof of Lemma 4.13.

Lemma 4.11. Let $A_{1}$ and $A_{3}$ be $k \times k$ real matrices. Assume that both $A_{1}$ and $A_{1} A_{3}$ are symmetric and $\sigma\left(A_{3}\right) \subset(-\infty, 0)$. Then

$$
\begin{equation*}
\operatorname{sgn} A_{1}+\operatorname{sgn}\left(A_{1} A_{3}\right)=0 \tag{4.4}
\end{equation*}
$$

Proof. It is clear that $A_{3}$ is invertible. We prove Lemma 4.11 in the following two steps.

Step 1. We assume that $A_{1}$ is invertible and proceed by induction on $k \in \mathbb{N}$.
If $k=1$, then $A_{1}, A_{3} \in \mathbb{R}$ and (4.4) obviously holds. Now assume (4.4) holds for $1 \leq k \leq l$. If we can prove 4.4 for $k=l+1$, then by mathematical induction (4.4) holds for any $k \in \mathbb{N}$ and Lemma 4.11 is proved in the case $A_{1}$ is invertible.

In view of the real Jordan canonical form decomposition of $A_{3}$, we only need to prove (4.4) for $k=l+1$ in the following two cases.

Case 1. There is an invertible $(l+1) \times(l+1)$ real matrix such that $Q^{-1} A_{3} Q$ is the $(l+1)$-order Jordan form

$$
\left(\begin{array}{cccccc}
\lambda & 1 & 0 & \cdots & \cdots & 0 \\
0 & \lambda & 1 & \ddots & \mathbf{0} & \vdots \\
0 & 0 & \ddots & \ddots & \ddots & \vdots \\
\vdots & \ddots & \ddots & \ddots & \ddots & 0 \\
\vdots & \mathbf{0} & \ddots & \ddots & \ddots & 1 \\
0 & \cdots & \cdots & 0 & 0 & \lambda
\end{array}\right):=\tilde{A}_{3}
$$

with $\lambda<0$.
Denoting by $\widetilde{A}_{1}=Q^{\top} A_{1} Q$, we have

$$
\tilde{A}_{1} \tilde{A}_{3}=Q^{\top} A_{1} Q Q^{-1} A_{3} Q=Q^{\top} A_{1} A_{3} Q
$$

Hence both matrices $\widetilde{A}_{1}$ and $\widetilde{A}_{1} \tilde{A}_{3}$ are symmetric and

$$
\begin{equation*}
\operatorname{sgn} A_{1}+\operatorname{sgn}\left(A_{1} A_{3}\right)=\operatorname{sgn} \tilde{A}_{1}+\operatorname{sgn}\left(\tilde{A}_{1} \tilde{A}_{3}\right) \tag{4.5}
\end{equation*}
$$

Since $\widetilde{A}_{1}=\left(a_{i, j}\right)_{1 \leq i, j \leq l+1}$ and $\widetilde{A}_{1} \widetilde{A}_{3}=\left(c_{i, j}\right)_{1 \leq i, j \leq l+1}$ are symmetric, $a_{i, j}=$ $a_{j . i}$ and $c_{i, j}=c_{j, i}$ for $1 \leq i, j \leq l+1$.
Claim 4.12. $a_{i, j}=0$ for $i+j \leq l+1$ and $a_{i, j}=a_{l+1,1}$ for $i+j=l+2$ with $1 \leq i, j \leq l+1$.

Proof. For $2 \leq j \leq l+1$, since $c_{1, j}=c_{j, 1}$,

$$
\lambda a_{1, j}+a_{1, j-1}=\lambda a_{j, 1}=\lambda a_{1, j} .
$$

Thus

$$
\begin{equation*}
a_{1, j-1}=0, \quad 2 \leq j \leq l+1 . \tag{4.6}
\end{equation*}
$$

For $2 \leq i, j \leq l+1$, since $c_{i, j}=c_{j, i}$ we have

$$
\lambda a_{i, j}+a_{i, j-1}=\lambda a_{j, i}+a_{j, i-1}=\lambda a_{i, j}+a_{i-1, j}
$$

So

$$
\begin{equation*}
a_{i, j-1}=a_{i-1, j}, \quad 2 \leq i, j \leq l+1 \tag{4.7}
\end{equation*}
$$

By (4.6) and (4.7) we have

$$
\begin{gather*}
a_{i, j}=a_{i-1, j+1}=\cdots=a_{2, i+j-2}=a_{1, i+j-1}=0, \\
1 \leq i, j \quad \text { and } \quad i+j \leq l+1,  \tag{4.8}\\
a_{l+1,1}=a_{l, 2}=a_{l-1,3}=\cdots=a_{2, l}=a_{1, l+1} . \tag{4.9}
\end{gather*}
$$

Hence, by (4.8) and (4.9), Claim4.12 is proved.
By Claim 4.12, let $a=a_{1, l+1}$; then

$$
\begin{align*}
\tilde{A}_{1} & =\left(\begin{array}{ccccccc}
0 & 0 & 0 & 0 & 0 & 0 & a \\
0 & 0 & 0 & 0 & 0 & a & * \\
0 & 0 & 0 & 0 & \cdot & * & * \\
0 & 0 & 0 & \cdot & * & * & * \\
0 & 0 & \cdot & * & * & * & * \\
0 & a & * & * & * & * & * \\
a & * & * & * & * & * & *
\end{array}\right)  \tag{4.10}\\
\tilde{A}_{1} \tilde{A}_{3} & =\left(\begin{array}{ccccccc}
0 & 0 & 0 & 0 & 0 & 0 & \lambda a \\
0 & 0 & 0 & 0 & 0 & \lambda a & * \\
0 & 0 & 0 & 0 & \cdot & * & * \\
0 & 0 & 0 & \cdot & * & * & * \\
0 & 0 & \cdot & * & * & * & * \\
0 & \lambda a & * & * & * & * & * \\
\lambda a & * & * & * & * & * & *
\end{array}\right) .
\end{align*}
$$

It is easy to see that $\widetilde{A}_{1} \widetilde{A}_{3}$ is congruent to $\lambda \widetilde{A}_{1}$. Since $\lambda<0$,

$$
\begin{gather*}
\operatorname{sgn}\left(\widetilde{A}_{1} \widetilde{A}_{3}\right)=\operatorname{sgn}\left(\lambda \tilde{A}_{1}\right)=-\operatorname{sgn}\left(\widetilde{A}_{1}\right), \\
\operatorname{sgn}\left(\widetilde{A}_{1} \widetilde{A}_{3}\right)+\operatorname{sgn} \widetilde{A}_{1}=0 . \tag{4.11}
\end{gather*}
$$

(4.5) and (4.11) imply (4.4). Hence Step 1 is proved in Case 1.

Case 2. There exists an invertible $(l+1) \times(l+1)$ real matrix $Q$ such that $Q^{-1} A_{3} Q=\operatorname{diag}\left(A_{4}, A_{5}\right)$, where $A_{4}$ is a $k_{1} \times k_{1}$ real matrix with $\sigma\left(A_{4}\right) \subset$ $(-\infty, 0)$ and $A_{5}$ is a $k_{2}$-order Jordan form

$$
A_{5}=\left(\begin{array}{ccccc}
\lambda & 1 & 0 & \cdots & 0 \\
0 & \ddots & \ddots & \ddots & \vdots \\
\vdots & \ddots & \ddots & \ddots & 0 \\
\vdots & \mathbf{0} & \ddots & \ddots & 1 \\
0 & \cdots & \cdots & 0 & \lambda
\end{array}\right)
$$

with $\lambda<0,1 \leq k_{1}, k_{2} \leq l$, and $k_{1}+k_{2}=l+1$.
We still denote $\widetilde{A}_{1}=Q^{\top} A_{1} Q$; then

$$
\tilde{A}_{1} \tilde{A}_{3}=Q^{\top} A_{1} Q Q^{-1} A_{3} Q=Q^{\top} A_{1} A_{3} Q
$$

So both $\widetilde{A}_{1}$ and $\widetilde{A}_{1} \widetilde{A}_{3}$ are symmetric and

$$
\begin{equation*}
\operatorname{sgn} A_{1}+\operatorname{sgn}\left(A_{1} A_{3}\right)=\operatorname{sgn} \tilde{A}_{1}+\operatorname{sgn}\left(\widetilde{A}_{1} \widetilde{A}_{3}\right) \tag{4.12}
\end{equation*}
$$

Correspondingly, we can write $\tilde{A}_{1}$ in the block form decomposition

$$
\tilde{A}_{1}=\left(\begin{array}{ll}
E_{1} & E_{2} \\
E_{2}^{\top} & E_{4}
\end{array}\right),
$$

where $E_{1}$ is a $k_{1} \times k_{1}$ real symmetric matrix and $E_{4}$ is a $k_{2} \times k_{2}$ real symmetric matrix. Then

$$
\tilde{A}_{1} \tilde{A}_{3}=\left(\begin{array}{ll}
E_{1} A_{4} & E_{2} A_{5} \\
E_{2}^{\top} A_{4} & E_{4} A_{5}
\end{array}\right)
$$

is symmetric.
Subcase 1. $E_{4}$ is invertible.
In this case we have

$$
\begin{align*}
&\left(\begin{array}{cc}
I_{k_{1}} & -E_{2} E_{4}^{-1} \\
0 & I_{k_{2}}
\end{array}\right)\left(\begin{array}{cc}
E_{1} & E_{2} \\
E_{2}^{\top} & E_{4}
\end{array}\right)\left(\begin{array}{cc}
I_{k_{1}} & 0 \\
-E_{4}^{-1} E_{2}^{\top} & I_{k_{2}}
\end{array}\right)=  \tag{4.13}\\
&\left(\begin{array}{ccc}
E_{1}-E_{2} E_{4}^{-1} E_{2}^{\top} & 0 \\
0 & E_{4}
\end{array}\right)
\end{align*}
$$

and

$$
\begin{align*}
& \left(\begin{array}{cc}
I_{k_{1}}-E_{2} E_{4}^{-1} & \\
0 & I_{k_{2}}
\end{array}\right)\left(\begin{array}{cc}
E_{1} A_{4} & E_{2} A_{5} \\
E_{2}^{\top} A_{4} & E_{4} A_{5}
\end{array}\right)\left(\begin{array}{cc}
I_{k_{1}} & 0 \\
-E_{4}^{-1} E_{2}^{\top} & I_{k_{2}}
\end{array}\right) \\
& \quad=\left(\begin{array}{cc}
E_{1} A_{4}-E_{2} E_{4}^{-1} E_{2}^{\top} A_{4} & 0 \\
0 & E_{4} A_{5}
\end{array}\right)  \tag{4.14}\\
& \quad=\left(\begin{array}{cc}
\left(E_{1}-E_{2} E_{4}^{-1} E_{2}^{\top}\right) A_{4} & 0 \\
0 & E_{4} A_{5}
\end{array}\right) .
\end{align*}
$$

Since the matrices $\widetilde{A}_{1}$ and $\tilde{A}_{1} \widetilde{A}_{3}$ are symmetric and invertible, by (4.13) and 4.14), both $E_{1}-E_{2} E_{4}^{-1} E_{2}^{\top}$ and $\left(E_{1}-E_{2} E_{4}^{-1} E_{2}^{\top}\right) A_{4}$ are symmetric and invertible. Hence from $1 \leq k_{1} \leq l, \sigma\left(A_{4}\right) \subset(-\infty, 0)$, and our induction hypothesis we obtain

$$
\begin{equation*}
\operatorname{sgn}\left(\left(E_{1}-E_{2} E_{4}^{-1} E_{2}^{\top}\right) A_{4}\right)+\operatorname{sgn}\left(E_{1}-E_{2} E_{4}^{-1} E_{2}^{\top}\right)=0 \tag{4.15}
\end{equation*}
$$

By (4.14), $E_{4} A_{5}$ is symmetric. Since $E_{4}$ is symmetric and invertible, $\sigma\left(A_{5}\right) \subset$ $(-\infty, 0)$ and $1 \leq k_{2} \leq l$, by our induction hypothesis we have

$$
\begin{equation*}
\operatorname{sgn}\left(E_{4} A_{5}\right)+\operatorname{sgn} E_{4}=0 \tag{4.16}
\end{equation*}
$$

From (4.13) we obtain

$$
\begin{equation*}
\operatorname{sgn} \tilde{A}_{1}=\operatorname{sgn}\left(E_{1}-E_{2} E_{4}^{-1} E_{2}^{\top}\right)+\operatorname{sgn} E_{4} \tag{4.17}
\end{equation*}
$$

By (4.14) there holds

$$
\begin{equation*}
\operatorname{sgn}\left(\tilde{A}_{1} \tilde{A}_{3}\right)=\operatorname{sgn}\left(\left(E_{1}-E_{2} E_{4}^{-1} E_{2}^{\top}\right) A_{4}\right)+\operatorname{sgn}\left(E_{4} A_{5}\right) \tag{4.18}
\end{equation*}
$$

Then by (4.15)-4.18) we have

$$
\begin{equation*}
\operatorname{sgn}\left(\tilde{A}_{1} \tilde{A}_{3}\right)+\operatorname{sgn} \tilde{A}_{1}=0 \tag{4.19}
\end{equation*}
$$

Therefore, 4.12 and 4.19) imply 4.4.
SUBCASE 2. $E_{4}$ is not invertible.
In this case we define $k_{2}$-order real invertible matrix

$$
E_{0}=\left(\begin{array}{ccccccc}
0 & 0 & 0 & 0 & 0 & 0 & 1 \\
0 & 0 & 0 & 0 & 0 & 1 & 0 \\
0 & 0 & 0 & 0 & . & 0 & 0 \\
0 & 0 & 0 & \cdot & 0 & 0 & 0 \\
0 & 0 & . & 0 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 & 0 & 0 & 0 \\
1 & 0 & 0 & 0 & 0 & 0 & 0
\end{array}\right)
$$

Then it is easy to verify that $E_{0} A_{5}$ is symmetric and $E_{4}+\varepsilon E_{0}$ is invertible for $0<\varepsilon \ll 1$. Define

$$
A_{\varepsilon}=\left(\begin{array}{cc}
E_{1} & E_{2} \\
E_{2}^{\top} & E_{4}+\varepsilon E_{0}
\end{array}\right)
$$

Since $\widetilde{A}_{1}$ and $\widetilde{A}_{1} \widetilde{A}_{3}$ are invertible, we have that both $A_{\varepsilon}$ and $A_{\varepsilon} \widetilde{A}_{3}$ are symmetric and invertible. Thus

$$
\begin{equation*}
\operatorname{sgn} \widetilde{A}_{1}=\operatorname{sgn} A_{\varepsilon}, \quad \operatorname{sgn}\left(\widetilde{A}_{1} \tilde{A}_{3}\right)=\operatorname{sgn}\left(A_{\varepsilon} \tilde{A}_{3}\right) \quad \text { for } 0<\varepsilon \ll 1 . \tag{4.20}
\end{equation*}
$$

By the proof of Subcase 1, we have

$$
\begin{equation*}
\operatorname{sgn}\left(A_{\varepsilon} \widetilde{A}_{3}\right)+\operatorname{sgn} A_{\varepsilon}=0 \tag{4.21}
\end{equation*}
$$

So from (4.20) we obtain

$$
\begin{equation*}
\operatorname{sgn}\left(\tilde{A}_{1} \tilde{A}_{3}\right)+\operatorname{sgn} \tilde{A}_{1}=0 . \tag{4.22}
\end{equation*}
$$

Then (4.4) holds from (4.22).
So in Case 2 (4.4) holds for $k=l+1$. Hence in the case $A_{1}$ is invertible, Lemma 4.11 holds and Step 1 is finished.

Step 2. We assume that $A_{1}$ is not invertible.
If $A_{1}=0$, (4.4) obviously holds.
If $1 \leq \operatorname{rank} A_{1}=m \leq k-1$, there is a real orthogonal matrix $G$ such that

$$
G^{\top} A_{1} G=\left(\begin{array}{cc}
0 & 0  \tag{4.23}\\
0 & \hat{A}_{1}
\end{array}\right),
$$

where $\hat{A}_{1}$ is an $m^{\text {th }}$-order invertible real symmetric matrix. Correspondingly, we write

$$
G^{-1} A_{3} G=\left(\begin{array}{ll}
F_{1} & F_{2} \\
F_{3} & F_{4}
\end{array}\right),
$$

where $F_{1}$ is a $(k-m) \times(k-m)$ real matrix and $F_{4}$ is a $m \times m$ real matrix.
Since $A_{1} A_{3}$ is symmetric, from

$$
G^{\top} A_{1} A_{3} G=G^{\top} A_{1} G G^{-1} A_{3} G=\left(\begin{array}{cc}
0 & 0 \\
\hat{A}_{1} F_{3} & \hat{A}_{1} F_{4}
\end{array}\right),
$$

we get $\widehat{A_{1}} F_{3}=0$. Hence $F_{3}=0$ by the invertibility of $\widehat{A_{1}}$. Therefore can write

$$
G^{-1} A_{3} G=\left(\begin{array}{cc}
F_{1} & F_{2}  \tag{4.24}\\
0 & F_{4}
\end{array}\right) .
$$

Hence

$$
G^{\top} A_{1} A_{3} G^{\top}=\left(\begin{array}{cc}
0 & 0  \tag{4.25}\\
0 & \hat{A}_{1} F_{4}
\end{array}\right)
$$

where $\hat{A}_{1} F_{4}$ is symmetric. Also, by (4.24) the matrix $F_{4}$ is invertible and $\sigma\left(F_{4}\right) \subset$ $(-\infty, 0)$. Thus by the proof of Step 1 , there holds

$$
\begin{equation*}
\operatorname{sgn}\left(\widehat{A_{1}} F_{4}\right)+\operatorname{sgn} \widehat{A}_{1}=0 \tag{4.26}
\end{equation*}
$$

Identities (4.23) and (4.25) give

$$
\begin{equation*}
\operatorname{sgn}\left(A_{1} A_{3}\right)+\operatorname{sgn} A_{1}=\operatorname{sgn}\left(\hat{A}_{1} F_{4}\right)+\operatorname{sgn} \widehat{A}_{1} . \tag{4.27}
\end{equation*}
$$

Then (4.26) and (4.27) give (4.4). Hence Step 2 is proved.

By Step 1 and Step 2 Lemma 4.11 holds.
We recall that the elliptic height $e(P)$ of $P$ is the total algebraic multiplicity of all eigenvalues of $P$ on $\mathbf{U}$ for any $P \in \operatorname{Sp}(2 n)$ (cf. [20, def. 1.8.1]).

LEMMA 4.13. Let

$$
R=\left(\begin{array}{ll}
A_{1} & I_{k} \\
A_{3} & A_{2}
\end{array}\right) \in \operatorname{Sp}(2 k)
$$

with $A_{3}$ being invertible. If $e\left(N_{k} R^{-1} N_{k} R\right)=2 m$, where $0 \leq m \leq k$, then

$$
\begin{equation*}
m-k \leq \frac{1}{2} \operatorname{sgn} M_{\varepsilon}(R) \leq k-m, \quad 0 \leq|\varepsilon| \ll 1 \tag{4.28}
\end{equation*}
$$

Proof. Since $e\left(N_{k} R^{-1} N_{k} R\right)=2 m$, there exists a symplectic matrix $P \in$ $\operatorname{Sp}(2 k)$ such that

$$
\begin{equation*}
P^{-1}\left(N_{k} R^{-1} N_{k} R\right) P=Q_{1} \diamond Q_{2} \tag{4.29}
\end{equation*}
$$

with $\sigma\left(Q_{1}\right) \in \mathbf{U}, \sigma\left(Q_{2}\right) \cap \mathbf{U}=\varnothing, Q_{1} \in \operatorname{Sp}(2 m)$, and $Q_{2} \in \operatorname{Sp}(2 k-2 m)$. By (ii) of Lemma 4.8, since $A_{3}$ is invertible we only need to prove 4.28) for $\varepsilon=0$.

Step 1. Assume $A_{1}$ is invertible.
Since $R$ is symplectic, we conclude from $R^{\top} J_{k} R=J_{k}$ that $A_{1}^{\top} A_{3}$ and $A_{2}$ are symmetric and

$$
A_{1}^{\top} A_{2}-A_{3}^{\top}=I_{k}
$$

Because $R^{\top}$ is also symplectic, $A_{1}$ is symmetric. Hence $A_{1} A_{3}$ is symmetric and

$$
\begin{equation*}
A_{1} A_{2}-A_{3}^{\top}=I_{k} \tag{4.30}
\end{equation*}
$$

By definition we have

$$
\begin{align*}
M_{0}(R) & =R^{\top}\left(\begin{array}{cc}
0 & -I_{k} \\
-I_{k} & 0
\end{array}\right) R+\left(\begin{array}{cc}
0 & I_{k} \\
I_{k} & 0
\end{array}\right) \\
& =-2\left(\begin{array}{cc}
A_{1} A_{3} & A_{3}^{\top} \\
A_{3} & A_{2}
\end{array}\right) . \tag{4.31}
\end{align*}
$$

Since $A_{1}$ is invertible, there holds

$$
\begin{align*}
& \left(\begin{array}{cc}
I_{k} & 0 \\
-A_{1}^{-1} & I_{k}
\end{array}\right)\left(\begin{array}{cc}
A_{1} A_{3} & A_{3}^{\top} \\
A_{3} & A_{2}
\end{array}\right)\left(\begin{array}{cc}
I_{k} & -A_{1}^{-1} \\
0 & I_{k}
\end{array}\right) \\
& \quad=\left(\begin{array}{cc}
A_{1} A_{3} & 0 \\
0 & -A_{1}^{-1} A_{3}^{\top}+A_{2}
\end{array}\right) \\
& \quad=\left(\begin{array}{cc}
A_{1} A_{3} & 0 \\
0 & A_{1}^{-1}
\end{array}\right) \tag{4.32}
\end{align*}
$$

where in the last equality we have used the equality (4.30). From (4.32) we obtain

$$
\frac{1}{2} \operatorname{sgn} M_{0}(R)=-\frac{1}{2} \operatorname{sgn}\left(\begin{array}{cc}
A_{1} A_{3} & 0  \tag{4.33}\\
0 & A_{1}^{-1}
\end{array}\right)
$$

By the Jordan canonical form decomposition of a complex matrix, there exists a complex invertible $k$-order matrix $G_{1}$ such that

$$
G_{1}^{-1} A_{3} G_{1}=\left(\begin{array}{ccccc}
u_{1} & * & * & * & * \\
0 & u_{2} & * & * & * \\
\vdots & \ddots & \ddots & * & * \\
\vdots & \mathbf{0} & \ddots & u_{k-1} & * \\
0 & \cdots & \cdots & 0 & u_{k}
\end{array}\right)
$$

with $u_{1}, u_{2}, \ldots, u_{k} \in \mathbb{C}$.
(4.2) gives

$$
N_{k} R^{-1} N_{k} R=I_{2 k}+2\left(\begin{array}{cc}
A_{3} & A_{2}  \tag{4.34}\\
A_{1} A_{3} & A_{3}^{\top}
\end{array}\right) .
$$

Since

$$
\left(\begin{array}{cc}
I_{k} & 0 \\
-A_{1} & I_{k}
\end{array}\right)\left(\begin{array}{cc}
A_{3} & A_{2} \\
A_{1} A_{3} & A_{3}^{\top}
\end{array}\right)\left(\begin{array}{cc}
I_{k} & 0 \\
A_{1} & I_{k}
\end{array}\right)=\left(\begin{array}{cc}
I_{k}+2 A_{3} & A_{2} \\
-A_{1} & -I_{k}
\end{array}\right),
$$

by (4.34) we have

$$
\left(\begin{array}{cc}
I_{k} & 0  \tag{4.35}\\
A_{1} & I_{k}
\end{array}\right)^{-1}\left(N_{k} R^{-1} N_{k} R\right)\left(\begin{array}{cc}
I_{k} & 0 \\
A_{1} & I_{k}
\end{array}\right)=\left(\begin{array}{cc}
3 I_{k}+4 A_{3} & 2 A_{2} \\
-2 A_{1} & -I_{k}
\end{array}\right):=R_{1} .
$$

By (4.35), for any $\lambda \in \mathbb{C}$ we get

$$
\lambda I_{2 k}-R_{1}=\left(\begin{array}{cc}
(\lambda-3) I_{k}-4 A_{3} & -2 A_{2}  \tag{4.36}\\
2 A_{1} & (\lambda+1) I_{k}
\end{array}\right) .
$$

Since $A_{1}$ is invertible, by 4.30) there holds

$$
\begin{align*}
& \left(\begin{array}{ccc}
I_{k} & -\frac{1}{2}\left((\lambda-3) I_{k}-4 A_{3}\right) A_{1}^{-1} \\
0 & I_{k}
\end{array}\right)\left(\begin{array}{cc}
(\lambda-3) I_{k}-4 A_{3} & -2 A_{2} \\
2 A_{1} & (\lambda+1) I_{k}
\end{array}\right)  \tag{4.37}\\
& =\left(\begin{array}{cc}
0 & -\frac{1}{2}\left(\left(\lambda^{2}-2 \lambda+1\right) I_{k}-4 \lambda A_{3}\right) A_{1}^{-1} \\
2 A_{1} & (\lambda+1) I_{k}
\end{array}\right) .
\end{align*}
$$

Then by (4.36)-(4.37) we have

$$
\begin{equation*}
\operatorname{det}\left(\lambda I_{2 k}-R_{1}\right)=\operatorname{det}\left(\left(\lambda^{2}-2 \lambda+1\right) I_{k}-4 \lambda A_{3}\right) \tag{4.38}
\end{equation*}
$$

Denote by $u_{1}, u_{2}, \ldots, u_{k}$ the $k$ complex eigenvalues of $A_{3}$; 4.38) gives

$$
\begin{align*}
\operatorname{det}\left(\lambda I_{2 k}-R_{1}\right) & =\prod_{i=1}^{k}\left(\lambda^{2}-2 \lambda+1-4 \lambda u_{i}\right) \\
& =\prod_{i=1}^{k}\left(\lambda^{2}-\left(2+4 u_{i}\right) \lambda+1\right) . \tag{4.39}
\end{align*}
$$

Thus from (4.35) and (4.39) we get

$$
\begin{align*}
\operatorname{det}\left(\lambda I_{2 k}-N_{k} R^{-1} N_{k} R\right) & =\prod_{i=1}^{k}\left(\lambda^{2}-2 \lambda+1-4 \lambda u_{i}\right) \\
& =\prod_{i=1}^{k}\left(\lambda^{2}-\left(2+4 u_{i}\right) \lambda+1\right) \tag{4.40}
\end{align*}
$$

It is easy to check that the equation $\lambda^{2}-\left(2+u_{i}\right) \lambda+1=0$ has two solutions on $\mathbf{U}$ if and only if $-4 \leq u_{i} \leq 0$ for $i=1, \ldots, k$. So by (4.29) without loss of generality we assume $u_{j} \in[-4,0)$ for $1 \leq j \leq m$ and $u_{j} \notin[-4,0)$ for $m+1 \leq j \leq k$. Then there exists a real invertible matrix of $k$-order $Q$ such that

$$
Q^{-1} A_{3} Q=\left(\begin{array}{cc}
A_{4} & 0 \\
0 & A_{5}
\end{array}\right):=\tilde{A}_{3}
$$

and $\sigma\left(A_{4}\right) \subset[-4,0), \sigma\left(A_{5}\right) \cap[-4,0)=\varnothing$, where $A_{4}$ is an $m$-order real invertible matrix and $A_{5}$ is a $(k-m)$-order real matrix.

Denote $\tilde{A}_{1}=Q^{\top} A_{1} Q$. We have

$$
\tilde{A}_{1} \tilde{A}_{3}=Q^{\top} A_{1} Q Q^{-1} A_{3} Q=Q^{\top} A_{1} A_{3} Q
$$

Hence both $\tilde{A}_{1}$ and $\tilde{A}_{1} \tilde{A}_{3}$ are symmetric, and we conclude that

$$
\begin{equation*}
\operatorname{sgn} A_{1}+\operatorname{sgn}\left(A_{1} A_{3}\right)=\operatorname{sgn} \tilde{A}_{1}+\operatorname{sgn}\left(\tilde{A}_{1} \tilde{A}_{3}\right) \tag{4.41}
\end{equation*}
$$

Correspondingly, we can write $\tilde{A_{1}}$ in the block form decomposition

$$
\tilde{A}_{1}=\left(\begin{array}{ll}
E_{1} & E_{2} \\
E_{2}^{\top} & E_{4}
\end{array}\right)
$$

where $E_{1}$ is an $m$-order real symmetric matrix and $E_{4}$ is a $(k-m)$-order real symmetric matrix. Then

$$
\tilde{A}_{1} \tilde{A}_{3}=\left(\begin{array}{ll}
E_{1} A_{4} & E_{2} A_{5} \\
E_{2}^{\top} A_{4} & E_{4} A_{5}
\end{array}\right)
$$

is symmetric.
By the same argument used in the proof of Subcase 2 of Lemma 4.11, without loss of generality we can assume $E_{1}$ is invertible (otherwise we can perturb it slightly so that it is invertible). So as in Subcase 1 of the proof of Lemma 4.11, we obtain

$$
\begin{array}{r}
\left(\begin{array}{cc}
I_{m} & 0 \\
-E_{2}^{\top} E_{1}^{-1} & I_{k-m}
\end{array}\right)\left(\begin{array}{cc}
E_{1} & E_{2} \\
E_{2}^{\top} & E_{4}
\end{array}\right)\left(\begin{array}{cc}
I_{m} & -E_{1}^{-1} E_{2} \\
0 & I_{k-m}
\end{array}\right)=  \tag{4.42}\\
\left(\begin{array}{cc}
E_{1} & 0 \\
0 & E_{4}-E_{2}^{\top} E_{1}^{-1} E_{2}
\end{array}\right)
\end{array}
$$

and

$$
\begin{gather*}
\left(\begin{array}{cc}
I_{m} & 0 \\
-E_{2}^{\top} E_{1}^{-1} & I_{k-m}
\end{array}\right)\left(\begin{array}{cc}
E_{1} A_{4} & E_{2} A_{5} \\
E_{2}^{\top} A_{4} & E_{4} A_{5}
\end{array}\right)\left(\begin{array}{cc}
I_{m} & -E_{1}^{-1} E_{2} \\
0 & I_{k-m}
\end{array}\right)=  \tag{4.43}\\
\left(\begin{array}{cc}
E_{1} A_{4} & 0 \\
0 & \left(E_{4}-E_{2}^{\top} E_{1}^{-1} E_{2}\right) A_{5}
\end{array}\right)
\end{gather*}
$$

By (4.43) we also have $E_{1} A_{4}$ is symmetric. Since $E_{1}$ is symmetric and invertible, $\sigma\left(A_{4}\right) \subset[-4,0)$, by Lemma 4.11$]$ we have

$$
\begin{equation*}
\operatorname{sgn}\left(E_{1} A_{4}\right)+\operatorname{sgn} E_{1}=0 \tag{4.44}
\end{equation*}
$$

By (4.42) and (4.42), there hold

$$
\begin{align*}
\operatorname{sgn} \tilde{A}_{1} & =\operatorname{sgn}\left(E_{4}-E_{2}^{\top} E_{1}^{-1} E_{2}\right)+\operatorname{sgn} E_{1},  \tag{4.45}\\
\operatorname{sgn}\left(\widetilde{A}_{1} \widetilde{A}_{3}\right) & =\operatorname{sgn}\left(\left(E_{4}-E_{2}^{\top} E_{1}^{-1} E_{2}\right) A_{5}\right)+\operatorname{sgn}\left(E_{1} A_{4}\right) . \tag{4.46}
\end{align*}
$$

(4.44)-(4.46) give

$$
\begin{align*}
\operatorname{sgn}\left(\tilde{A}_{1} \tilde{A}_{3}\right)+\operatorname{sgn} \tilde{A}_{1}= & \operatorname{sgn}\left(\left(E_{4}-E_{2}^{\top} E_{1}^{-1} E_{2}\right) A_{5}\right) \\
& +\operatorname{sgn}\left(E_{4}-E_{2}^{\top} E_{1}^{-1} E_{2}\right)  \tag{4.47}\\
\in & {[-2(k-m), 2(k-m)] . }
\end{align*}
$$

Then (4.28) holds from (4.33), 4.41), and 4.47).
Step 2. Assume $A_{1}$ is not invertible.
If $A_{1}=0$, then $A_{3}=-I_{k}$ and $m=k$. It is easy to check that

$$
M_{0}(R)=2\left(\begin{array}{cc}
0 & I_{k} \\
I_{k} & -A_{2}
\end{array}\right) \quad \text { is congruent to } \quad 2\left(\begin{array}{cc}
0 & I_{k} \\
I_{k} & 0
\end{array}\right),
$$

so $\operatorname{sgn} M_{0}(R)=0$ and (4.28) holds.
If $1 \leq \operatorname{rank} A_{1}=r \leq k-1$, there is a $k \times k$ invertible matrix $G$ with $\operatorname{det} G>0$ such that

$$
\begin{equation*}
\left(G^{-1}\right)^{\top} A_{1} G^{-1}=\operatorname{diag}(0, \Lambda) \tag{4.48}
\end{equation*}
$$

where $\Lambda$ is a $r \times r$ real invertible matrix. Hence

$$
\begin{align*}
\operatorname{diag}\left(\left(G^{\top}\right)^{-1}, G\right) \cdot R \cdot \operatorname{diag}\left(G^{-1}, G^{\top}\right) & =\left(\begin{array}{cc}
\left(G^{\top}\right)^{-1} A_{1} G^{-1} & I_{k} \\
G A_{3} G^{-1} & G A_{2} G^{\top}
\end{array}\right) \\
\text { 49) } \quad & :=R_{2}=\left(\begin{array}{cccc}
0 & 0 & I_{k-r} & 0 \\
0 & \Lambda & 0 & I_{r} \\
B_{1} & B_{2} & D_{1} & D_{2} \\
B_{3} & B_{4} & D_{3} & D_{4}
\end{array}\right), \tag{4.49}
\end{align*}
$$

where $B_{1}$ and $D_{1}$ are $(k-r) \times(k-r)$ matrices and $B_{4}$ and $D_{4}$ are $r \times r$ matrices.

Since $R_{2}$ is symplectic and $\Lambda$ is invertible, there holds $R_{2}^{\top} J_{k} R_{2}=J_{k}$. It implies that $B_{3}=0, D_{3}=D_{2}^{\top}, B_{1}=-I_{k-r}$, and $D_{1}$ and $D_{4}$ are symmetric. Thus

$$
R_{2}=\left(\begin{array}{cccc}
0 & 0 & I_{k-r} & 0 \\
0 & \Lambda & 0 & I_{r} \\
B_{1} & B_{2} & D_{1} & D_{2} \\
0 & B_{4} & D_{2}^{\top} & D_{4}
\end{array}\right)
$$

For $t \in[0,1]$, we define

$$
\beta(t)=\left(\begin{array}{cccc}
0 & 0 & I_{k-r} & 0 \\
0 & \Lambda & 0 & I_{r} \\
B_{1} & t B_{2} & t D_{1} & t D_{2} \\
0 & B_{4} & t D_{2}^{\top} & D_{4}
\end{array}\right)
$$

It is easy to check that $\beta$ is a symplectic path and $\nu_{L_{j}}(\beta(t))=0$ for all $t \in[0,1]$ and $j=0,1$. We also have $\beta(1)=R_{2}$ and

$$
\beta(0)=\left(\begin{array}{cccc}
0 & 0 & I_{k-r} & 0 \\
0 & \Lambda & 0 & I_{r} \\
B_{1} & 0 & 0 & 0 \\
0 & B_{4} & 0 & D_{4}
\end{array}\right)=-J_{k-r} \diamond\left(\begin{array}{cc}
\Lambda & I_{r} \\
B_{4} & D_{4}
\end{array}\right):=R_{3}
$$

Then by lemma 2.2 of [32], Lemma 4.9, and Remark 4.3] we have

$$
\begin{align*}
\frac{1}{2} \operatorname{sgn} M_{0}\left(R_{2}\right) & =\frac{1}{2} \operatorname{sgn} M_{0}\left(-J_{k-r}\right)+\frac{1}{2} \operatorname{sgn} M_{0}\left(\left(\begin{array}{cc}
\Lambda & I_{r} \\
B_{4} & D_{4}
\end{array}\right)\right) \\
& =\frac{1}{2} \operatorname{sgn} M_{0}\left(\left(\begin{array}{cc}
\Lambda & I_{r} \\
B_{4} & D_{4}
\end{array}\right)\right) \tag{4.50}
\end{align*}
$$

Since $R_{2} \sim R$, by (4.50) we have

$$
\frac{1}{2} \operatorname{sgn} M_{0}(R)=\frac{1}{2} \operatorname{sgn} M_{0}\left(\left(\begin{array}{cc}
\Lambda & I_{r}  \tag{4.51}\\
B_{4} & D_{4}
\end{array}\right)\right)
$$

By (4.2), there holds

$$
N_{k} R_{2}^{-1} N_{k} R_{2}=I_{2 k}+2\left(\begin{array}{cccc}
B_{1} & B_{2} & D_{1} & D_{2}  \tag{4.52}\\
0 & B_{4} & D_{2}^{\top} & D_{4} \\
0 & 0 & B_{1}^{\top} & 0 \\
0 & \Lambda B_{4} & B_{2}^{\top} & B_{4}^{\top}
\end{array}\right)
$$

By 4.52 for any $\lambda \in \mathbb{C}$, we obtain

$$
\begin{align*}
& \operatorname{det}\left(\lambda I_{2 k}-N_{k} R_{2}^{-1} N_{k} R_{2}\right) \\
& =\operatorname{det}\left((\lambda-1) I_{k-r}-2 B_{1}\right) \operatorname{det}\left((\lambda-1) I_{k-r}-2 B_{1}^{\top}\right) \\
& \quad \cdot \operatorname{det}\left(\begin{array}{cc}
(\lambda-1) I_{r}-2 B_{4} & -2 D_{4} \\
-2 \Lambda B_{4} & (\lambda-1) I_{r}-2 B_{4}^{\top}
\end{array}\right) \\
& =\operatorname{det}\left(\lambda I_{2 k}-N_{k} R_{3}^{-1} N_{k} R_{3}\right) \tag{4.53}
\end{align*}
$$

where

$$
N_{k} R_{3}^{-1} N_{k} R_{3}=I_{2 k}+2\left(\begin{array}{cccc}
B_{1} & 0 & 0 & 0 \\
0 & B_{4} & 0 & D_{4} \\
0 & 0 & B_{1}^{\top} & 0 \\
0 & \Lambda B_{4} & 0 & B_{4}^{\top}
\end{array}\right) .
$$

So (4.53) gives

$$
\begin{equation*}
\sigma\left(N_{k} R^{-1} N_{k} R\right)=\sigma\left(N_{k} R_{2}^{-1} N_{k} R_{2}\right)=\sigma\left(N_{k} R_{3}^{-1} N_{k} R_{3}\right) . \tag{4.54}
\end{equation*}
$$

Since $B_{1}=-I_{k-r}$ and

$$
R_{3}=\left(-J_{k-r}\right) \diamond\left(\begin{array}{cc}
\Lambda & I_{r} \\
B_{4} & D_{4}
\end{array}\right),
$$

(4.54) gives

$$
e\left(N_{r}\left(\begin{array}{cc}
\Lambda & I_{r}  \tag{4.55}\\
B_{4} & D_{4}
\end{array}\right)^{-1} N_{r}\left(\begin{array}{cc}
\Lambda & I_{r} \\
B_{4} & D_{4}
\end{array}\right)\right)=2(m-(k-r))
$$

Step 1 implies that

$$
\frac{1}{2}\left|\operatorname{sgn} M_{0}\left(\left(\begin{array}{cc}
\Lambda & I_{r}  \tag{4.56}\\
B_{4} & D_{4}
\end{array}\right)\right)\right| \leq r-(m-(k-r))=k-m .
$$

Then (4.28) follows from (4.51) and (4.56). This finishes the proof of Step 2.
With Step 1 and Step 2, the proof of Lemma 4.13 is completed.
The following result is about the ( $L_{0}, L_{1}$ )-normal forms of $L_{0}$-degenerate symplectic matrices, which generalizes lemma 2.10 of [33].
Lemma 4.14. Let $R \in \operatorname{Sp}(2 k)$ have the block form

$$
R=\left(\begin{array}{ll}
A & B \\
C & D
\end{array}\right) \quad \text { with } 1 \leq \operatorname{rank} B=r<k
$$

We have
(i)

$$
R \sim\left(\begin{array}{cccc}
A_{1} & B_{1} & I_{r} & 0 \\
0 & D_{1} & 0 & 0 \\
A_{3} & B_{3} & A_{2} & 0 \\
C_{3} & D_{3} & C_{2} & D_{2}
\end{array}\right)
$$

where $A_{1}, A_{2}, A_{3}$ are $r \times r$ matrices, $D_{1}, D_{2}, D_{3}$ are $(k-r) \times(k-r)$ matrices, $B_{1}, B_{3}$ are $r \times(k-r)$ matrices, and $C_{2}, C_{3}$ are $(k-r) \times r$ matrices.
(ii) If $A_{3}$ is invertible, we have

$$
R \sim\left(\begin{array}{cc}
A_{1} & I_{r}  \tag{4.57}\\
A_{3} & A_{2}
\end{array}\right) \diamond\left(\begin{array}{cc}
D_{1} & 0 \\
\tilde{D}_{3} & D_{2}
\end{array}\right)
$$

where $\widetilde{D}_{3}$ is a $(k-r) \times(k-r)$ matrix.
(iii) If $1 \leq \operatorname{rank} A_{3}=\lambda \leq r-1$, then

$$
R \sim\left(\begin{array}{cc}
U & I_{\lambda}  \tag{4.58}\\
\Lambda & V
\end{array}\right) \diamond\left(\begin{array}{cccc}
\tilde{A}_{1} & \widetilde{B}_{1} & I_{r-\lambda} & 0 \\
0 & D_{1} & 0 & 0 \\
0 & \widetilde{B}_{3} & \tilde{A}_{2} & 0 \\
\widetilde{C}_{3} & \widetilde{D}_{3} & \widetilde{C}_{2} & \widetilde{D}_{2}
\end{array}\right),
$$

where $\tilde{A}_{1}, \tilde{A}_{2}$ are $(r-\lambda) \times(r-\lambda)$ matrices, $\widetilde{B}_{1}, \widetilde{B}_{3} \underset{\sim}{\operatorname{Dre}} \underset{\sim}{r}(r-\lambda) \times(k-r)$ matrices, $\widetilde{C}_{2}, \widetilde{C}_{3}$ are $(k-r) \times(r-\lambda)$ matrices, $D_{1}, \widetilde{D}_{2}, \widetilde{D}_{3}$ are $(k-r) \times$ $(k-r)$ matrices, $U, V, \Lambda$ are $\lambda \times \lambda$ matrices, and $\Lambda$ is invertible.
(iv) If $A_{3}=0$, then $A_{1}, A_{2}$ are symmetric and $A_{1} A_{2}=I_{r}$. Suppose $m^{+}\left(A_{1}\right)=$ $p, m^{-}\left(A_{1}\right)=r-p$, and $0 \leq \operatorname{rank} B_{3}=\lambda \leq \min \{r, k-r\}$, then

$$
\begin{align*}
& N_{k} R^{-1} N_{k} R \approx\left(\begin{array}{ll}
1 & 1 \\
0 & 1
\end{array}\right)^{\diamond p+q^{-}} \diamond\left(\begin{array}{rr}
1 & -1 \\
0 & 1
\end{array}\right)^{\diamond\left(r-p+q^{+}\right)} \diamond I_{2}^{\diamond q^{0}} \diamond D(2)^{\diamond \lambda},  \tag{4.59}\\
& m^{+}\left(A^{\top} C\right)=\lambda+q^{+},  \tag{4.60}\\
& m^{0}\left(A^{\top} C\right)=r-\lambda+q^{0},  \tag{4.61}\\
& m^{-}\left(A^{\top} C\right)=\lambda+q^{-}, \tag{4.62}
\end{align*}
$$

where $q^{*} \geq 0$ for $*= \pm, 0, q^{+}+q^{0}+q^{-}=k-r-\lambda, M^{\diamond 0}$ means the corresponding component does not appear at all for $M$ being one of the four matrices on the right-hand side of (4.59).

Proof. By lemma 2.10 of [33] or the same argument used in the proof of theorem 3.1 of [34], (i)-(iii) hold. So we only need to prove 4.59-4.62).

By (i) and $A_{3}=0$ we have

$$
R \sim\left(\begin{array}{cccc}
A_{1} & B_{1} & I_{r} & 0  \tag{4.63}\\
0 & D_{1} & 0 & 0 \\
0 & B_{3} & A_{2} & 0 \\
C_{3} & D_{3} & C_{2} & D_{2}
\end{array}\right):=R_{1}
$$

Since $R_{1}$ is symplectic we have $R_{1}^{\top} J_{k} R_{1}=J_{k}$. Then we have $A_{1}, A_{2}$ are symmetric and $A_{1} A_{2}=I_{r} . D_{1} D_{2}^{\top}=I_{k-r}$ and $A_{1}^{\top} B_{3}=C_{3}^{\top} D_{1}$. 4.2) yields

$$
N_{k} R_{1}^{-1} N_{k} R_{1}=\left(\begin{array}{cccc}
I_{r} & 2 B_{3} & 2 A_{2} & 0  \tag{4.64}\\
0 & I_{k-r} & 0 & 0 \\
0 & 2 A_{1}^{\top} B_{3} & I_{r} & 0 \\
2 B_{3}^{\top} A_{1} & 2 B_{1}^{\top} B_{3}+2 D_{1}^{\top} D_{3} & 2 B_{3}^{\top} & I_{k-r}
\end{array}\right)
$$

By Remark 4.6 we obtain

$$
m^{*}\left(A^{\top} C\right)=m^{*}\left(\left(\begin{array}{cc}
0 & A_{1}^{\top} B_{3}  \tag{4.65}\\
B_{3}^{\top} A_{1} & B_{1}^{\top} B_{3}+D_{1}^{\top} D_{3}
\end{array}\right)\right), \quad *=+,-, 0
$$

Since $0 \leq \operatorname{rank} B_{3}=\lambda \leq \min \{r, k-r\}$, there exist $r \times r$ and $(k-r) \times(k-r)$ real invertible matrices $G_{1}$ and $G_{2}$ such that

$$
G_{1} B_{3} G_{2}=\left(\begin{array}{cc}
I_{\lambda} & 0  \tag{4.66}\\
0 & 0
\end{array}\right):=F
$$

Note that if $\lambda=0$ then $B_{3}=01$ if $\lambda=\min \{r, k-r\}$ then

$$
G_{1} B_{3} G_{2}=\left(\begin{array}{ll}
I_{\lambda} & 0
\end{array}\right) \quad \text { or } \quad\binom{c I_{\lambda}}{0}
$$

if $\lambda=r=k-r$ then $G_{1} B_{3} G_{2}=I_{\lambda}$. The proof below can still go through by a suitable adjustment.

By (4.66 we have

$$
\begin{align*}
& \left(\begin{array}{cc}
G_{1} A_{1}^{-1} & 0 \\
0 & G_{2}^{\top}
\end{array}\right)\left(\begin{array}{cc}
0 & A_{1}^{\top} B_{3} \\
B_{3}^{\top} A_{1} & B_{1}^{\top} B_{3}+D_{1}^{\top} D_{3}
\end{array}\right)\left(\begin{array}{cc}
A_{1}^{-1} G_{1}^{\top} & 0 \\
0 & G_{2}
\end{array}\right) \\
& \quad=\left(\begin{array}{cc}
0 & G_{1} B_{3} G_{2} \\
G_{2}^{\top} B_{3}^{\top} G_{1}^{\top} & U
\end{array}\right)=\left(\begin{array}{cccc}
0 & 0 & I_{\lambda} & 0 \\
0 & 0 & 0 & 0 \\
I_{\lambda} & 0 & U_{1} & U_{2} \\
0 & 0 & U_{2}^{\top} & U_{4}
\end{array}\right) \tag{4.67}
\end{align*}
$$

Then

$$
\begin{align*}
& \left(\begin{array}{cccc}
I_{\lambda} & 0 & 0 & 0 \\
0 & I_{r-\lambda} & 0 & 0 \\
-\frac{1}{2} U_{1} & 0 & I_{\lambda} & 0 \\
-U_{2}^{\top} & 0 & 0 & I_{k-r-\lambda}
\end{array}\right)\left(\begin{array}{cccc}
0 & 0 & I_{\lambda} & 0 \\
0 & 0 & 0 & 0 \\
I_{\lambda} & 0 & U_{1} & U_{2} \\
0 & 0 & U_{2}^{\top} & U_{4}
\end{array}\right) \\
& \left(\begin{array}{cccc}
I_{\lambda} & 0 & -\frac{1}{2} U_{1} & -U_{2} \\
0 & I_{r-\lambda} & 0 & 0 \\
0 & 0 & I_{\lambda} & 0 \\
0 & 0 & 0 & I_{k-r-\lambda}
\end{array}\right)  \tag{4.68}\\
& \quad=\left(\begin{array}{cccc}
0 & 0 & I_{\lambda} & 0 \\
0 & 0 & 0 & 0 \\
I_{\lambda} & 0 & 0 & 0 \\
0 & 0 & 0 & U_{4}
\end{array}\right)
\end{align*}
$$

Set

$$
\begin{equation*}
q^{*}=m^{*}\left(U_{4}\right), \quad *= \pm, 0 \tag{4.69}
\end{equation*}
$$

Then $q^{+}+q^{0}+q^{-}=k-r-\lambda$ and (4.60-4.62) hold from 4.65), 4.67), and (4.68).

Also by 4.68) and Lemma 4.4 we have

$$
\left(\begin{array}{cc}
I_{k-r-\lambda} & 0  \tag{4.70}\\
2 U_{4} & I_{k-r-\lambda}
\end{array}\right) \approx\left(\begin{array}{cc}
1 & 1 \\
0 & 1
\end{array}\right)^{\diamond q^{-}} \diamond I_{2}^{\diamond q^{0}} \diamond\left(\begin{array}{cc}
1 & -1 \\
0 & 1
\end{array}\right)^{\diamond q^{+}}
$$

By (4.67), there holds

$$
\begin{align*}
& \operatorname{diag}\left(\left(G_{1}^{\top}\right)^{-1} A_{1}, G_{2}^{-1}, G_{1} A_{1}^{-1}, G_{2}^{\top}\right)\left(N_{k} R_{1}^{-1} N_{k} R_{1}\right) \\
& \quad \cdot \operatorname{diag}\left(A_{1}^{-1} G_{1}^{\top}, G_{2}, A_{1} G_{1}^{-1},\left(G_{2}^{\top}\right)^{-1}\right) \\
& \quad=\left(\begin{array}{cccc}
I_{r} & 2 E & 2 \tilde{A}_{1} & 0 \\
0 & I_{k-r} & 0 & 0 \\
0 & 2 F & I_{r} & 0 \\
2 F^{\top} & 2 U & 2 E^{\top} & I_{k-r}
\end{array}\right):=M \tag{4.71}
\end{align*}
$$

where $\tilde{A}_{1}=\left(G_{1}^{\top}\right)^{-1} A_{1} G_{1}^{-1}$ and $E=\left(G_{1}^{\top}\right)^{-1} A_{1} B_{3} G_{2}=\tilde{A}_{1} F$.
Since $M$ is symplectic, we have $M^{\top} J_{k} M=J_{k}$. Then we have $E=\tilde{A_{1}} F$. Since $\widetilde{A}_{1}=\left(G_{1}^{\top}\right)^{-1} A_{1} G_{1}^{-1}$, it is congruent to $\operatorname{diag}\left(a_{1}, a_{2}, \ldots, a_{r}\right)$ with

$$
\begin{array}{ll}
a_{i}=1, & \\
a_{j}=-1 \leq i \leq p  \tag{4.72}\\
a_{j} & p+1 \leq j \leq r \text { for some } 0 \leq p \leq r
\end{array}
$$

Then there is an invertible $r \times r$ real matrix $Q$ such that $\operatorname{det} Q>0$ and

$$
\begin{align*}
Q \tilde{A}_{1} Q^{\top} & =\operatorname{diag}\left(a_{1}, a_{2}, \ldots, a_{r}\right) \\
& =\operatorname{diag}\left(\operatorname{diag}\left(a_{1}, a_{2}, \ldots, a_{\lambda}\right), \operatorname{diag}\left(a_{\lambda+1}, \ldots, a_{r}\right)\right)  \tag{4.73}\\
& :=\operatorname{diag}\left(\Lambda_{1}, \Lambda_{2}\right)
\end{align*}
$$

Since $\operatorname{det} Q>0$ we can join it to $I_{r}$ by an invertible continuous matrix path. So there is a continuous invertible symmetric matrix path $\alpha_{1}$ such that $\alpha_{1}(1)=\widetilde{A}_{1}$ and $\alpha_{1}(0)=\operatorname{diag}\left(a_{1}, a_{2}, \ldots, a_{r}\right)$ with

$$
m^{*}\left(\alpha_{1}(t)\right)=m^{*}\left(\tilde{A}_{1}\right)=m^{*}\left(A_{1}\right), \quad t \in[0,1], *=+,-
$$

Define symmetric matrix path

$$
\alpha_{2}(t)=\left(\begin{array}{cc}
2 t U_{1} & 2 t U_{2} \\
2 t U_{2}^{\top} & 2 U_{4}
\end{array}\right), \quad t \in[0,1]
$$

For $t \in[0,1]$, define

$$
\beta(t)=\left(\begin{array}{cccc}
I_{r} & 2 \alpha_{1}(t) F & 2 \alpha_{1}(t) & 0 \\
0 & I_{k-r} & 0 & 0 \\
0 & 2 F & I_{r} & 0 \\
2 F^{\top} & \alpha_{2}(t) & 2 F^{\top} \alpha_{1}(t)^{\top} & I_{k-r}
\end{array}\right)
$$

Then since $M$ is symplectic, it is easy to check that $\beta$ is a continuous path of symplectic matrices. Since

$$
F=\left(\begin{array}{cc}
I_{\lambda} & 0 \\
0 & 0
\end{array}\right)
$$

and $\alpha_{1}(t)$ is invertible, by direct computation, we have

$$
\begin{aligned}
\operatorname{rank}\left(\beta(t)-I_{2 k}\right) & =2 \lambda+\operatorname{rank}\left(\alpha_{1}(t)\right)+\operatorname{rank}\left(U_{4}\right) \\
& =2 \lambda+r+m^{+}\left(U_{4}\right)+m^{-}\left(U_{4}\right)
\end{aligned}
$$

Hence

$$
v_{1}(\beta(t))=v_{1}(\beta(1))=v_{1}(M), \quad t \in[0,1]
$$

Because $\sigma(\beta(t))=\{1\}$, by Definition 2.2 and Lemma 2.4

$$
\begin{aligned}
M & =\beta(1) \approx \beta(0) \\
& =\left(\begin{array}{cccccc}
I_{\lambda} & 0 & 2 \Lambda_{1} & 2 \Lambda_{1} & 0 & 0 \\
0 & I_{r-\lambda} & 0 & 0 & 2 \Lambda_{2} & 0 \\
0 & 0 & I_{\lambda} & 0 & 0 & 0 \\
0 & 0 & 2 I_{\lambda} & I_{\lambda} & 0 & 0 \\
0 & 0 & 0 & 0 & I_{r-\lambda} & 0 \\
2 I_{\lambda} & 0 & 0 & 2 \Lambda_{1} & 0 & I_{\lambda}
\end{array}\right) \diamond\left(\begin{array}{cc}
I_{k-r}-\lambda & 0 \\
2 U_{4} & I_{k-r-\lambda}
\end{array}\right) \\
& \approx\left(\begin{array}{cccc}
I_{\lambda} & 2 \Lambda_{1} & 2 \Lambda_{1} & 0 \\
0 & I_{\lambda} & 0 & 0 \\
0 & 2 I_{\lambda} & I_{\lambda} & 0 \\
2 I_{\lambda} & 0 & 2 \Lambda_{1} & I_{\lambda}
\end{array}\right) \diamond\left(\begin{array}{ccc}
I_{r-\lambda} & 2 \Lambda_{2} \\
0 & I_{r-\lambda}
\end{array}\right) \diamond\left(\begin{array}{cc}
I_{k-r-\lambda} & 0 \\
2 U_{4} & I_{k-r-\lambda}
\end{array}\right) \\
& =\left(\begin{array}{cccc}
I_{\lambda} & 2 \Lambda_{1} & 2 \Lambda_{1} & 0 \\
0 & I_{\lambda} & 0 & 0 \\
0 & 2 I_{\lambda} & I_{\lambda} & 0 \\
2 I_{\lambda} & 0 & 2 \Lambda_{1} & I_{\lambda}
\end{array}\right) \diamond \diamond_{j=\lambda+1}^{r}\left(\begin{array}{ll}
1 & 2 a_{j} \\
0 & 1
\end{array}\right) \diamond\left(\begin{array}{cc}
I_{k-r}-\lambda & 0 \\
2 U_{4} & I_{k-r}-\lambda
\end{array}\right)
\end{aligned}
$$

We define the continuous symplectic matrix path

$$
\psi(t)=\left(\begin{array}{cccc}
I_{\lambda} & 2\left(1-t^{2}\right) \Lambda_{1} & 2 \Lambda_{1} & 0 \\
0 & (1+t) I_{\lambda} & 0 & 0 \\
0 & 2\left(1-t^{2}\right) I_{\lambda} & I_{\lambda} & 0 \\
2(1-t) I_{\lambda} & 0 & 2(1-t) \Lambda_{1} & \frac{1}{1+t} I_{\lambda}
\end{array}\right), \quad t \in[0,1]
$$

Since $\Lambda_{1}$ is invertible, $v(\psi(t)) \equiv \lambda$ for $t \in[0,1]$. So by $\sigma(\psi(t)) \cap \mathbf{U}=\{1\}$ for $t \in[0, t]$ and Definition 2.2 we obtain

$$
\begin{align*}
\left(\begin{array}{cccc}
I_{\lambda} & \Lambda_{1} & 2 \Lambda_{1} & 0 \\
0 & I_{\lambda} & 0 & 0 \\
0 & 2 I_{\lambda} & I_{\lambda} & 0 \\
2 I_{\lambda} & 0 & 2 \Lambda_{1} & I_{\lambda}
\end{array}\right) & =\psi(0) \approx \psi(1) \\
& =\left(\begin{array}{cc}
I_{\lambda} & 2 \Lambda_{1} \\
0 & I_{\lambda}
\end{array}\right) \diamond\left(\begin{array}{cc}
2 I_{\lambda} & 0 \\
0 & \frac{1}{2} I_{\lambda}
\end{array}\right) \\
& =\diamond_{j=1}^{\lambda}\left(\begin{array}{cc}
1 & 2 a_{j} \\
0 & 1
\end{array}\right) \diamond D(2)^{\diamond \lambda} \tag{4.74}
\end{align*}
$$

Thus by (4.74), 4.74), and Remark 2.3 we get

$$
M \approx\left(\diamond_{j=1}^{r}\left(\begin{array}{cc}
1 & a_{j}  \tag{4.75}\\
0 & 1
\end{array}\right)\right) \diamond D(2)^{\diamond \lambda} \diamond\left(\begin{array}{cc}
I_{k-r-\lambda} & 0 \\
U_{4} & I_{k-r-\lambda}
\end{array}\right)
$$

So by (4.70), 4.72, and Remark 2.3, there holds

$$
M \approx\left(\begin{array}{ll}
1 & 1  \tag{4.76}\\
0 & 1
\end{array}\right)^{\diamond\left(p+q^{-}\right)} \diamond\left(\begin{array}{rr}
1 & -1 \\
0 & 1
\end{array}\right)^{\diamond\left(r-p+q^{+}\right)} \diamond I_{2}^{\diamond q^{0}} \diamond D(2)^{\diamond \lambda}
$$

By Lemma 4.7, (4.63), and (4.71), we have

$$
\begin{equation*}
N_{k} R^{-1} N_{k} R \approx M \tag{4.77}
\end{equation*}
$$

Then 4.59 holds from (4.76) and 4.77). The proof of Lemma 4.14 is completed.

## 5 The Mixed ( $L_{0}, L_{1}$ )-Concavity

DEFINITION 5.1. The mixed $\left(L_{0}, L_{1}\right)$-concavity and mixed $\left(L_{1}, L_{0}\right)$-concavity of a symplectic path $\gamma \in \mathcal{P}_{\tau}(2 n)$ are defined respectively by

$$
\mu_{\left(L_{0}, L_{1}\right)}(\gamma)=i_{L_{0}}(\gamma)-v_{L_{1}}(\gamma), \quad \mu_{\left(L_{1}, L_{0}\right)}(\gamma)=i_{L_{1}}(\gamma)-v_{L_{0}}(\gamma)
$$

Proposition C of [21], proposition 6.1 of [17], and Theorem 4.2 imply the following result:

## Proposition 5.2. There hold

$$
\begin{align*}
\mu_{\left(L_{0}, L_{1}\right)}(\gamma)+\mu_{\left(L_{1}, L_{0}\right)}(\gamma)= & i\left(\gamma^{2}\right)-v\left(\gamma^{2}\right)-n  \tag{5.1}\\
\mu_{\left(L_{0}, L_{1}\right)}(\gamma)-\mu_{\left(L_{1}, L_{0}\right)}(\gamma)= & \operatorname{concav}_{\left(L_{0}, L_{1}\right)}^{*}(\gamma)=\frac{1}{2} \operatorname{sgn} M_{\varepsilon}(\gamma(\tau)),  \tag{5.2}\\
& 0<-\varepsilon \ll 1
\end{align*}
$$

Theorem 1.10 in Section 1 is a special case of the following result:
THEOREM 5.3. For $\gamma \in \mathcal{P}_{\tau}(2 n)$, let $P=\gamma(\tau)$. If $i_{L_{0}}(\gamma) \geq 0, i_{L_{1}}(\gamma) \geq 0$, $i(\gamma) \geq n$, and $\gamma^{2}(t)=\gamma(t-\tau) \gamma(\tau)$ for all $t \in[\tau, 2 \tau]$, then

$$
\begin{align*}
& \mu_{\left(L_{0}, L_{1}\right)}(\gamma)+S_{P^{2}}^{+}(1) \geq 0  \tag{5.3}\\
& \mu_{\left(L_{1}, L_{0}\right)}(\gamma)+S_{P^{2}}^{+}(1) \geq 0 \tag{5.4}
\end{align*}
$$

Proof. The proofs of (5.3) and (5.4) are almost the same. We only prove (5.4), which yields Theorem 1.10 .

Claim 5.4. Under the conditions of Theorem5.3, if

$$
P^{2} \approx\left(\begin{array}{ll}
1 & 1  \tag{5.5}\\
0 & 1
\end{array}\right)^{\diamond p_{1}} \diamond D(2)^{\diamond p_{2}} \diamond \widetilde{P}
$$

then

$$
\begin{equation*}
i\left(\gamma^{2}\right)+2 S_{P^{2}}^{+}(1)-v\left(\gamma^{2}\right) \geq n+p_{1}+p_{2} \tag{5.6}
\end{equation*}
$$

Proof of Claim 5.4. By theorem 7.8 of [19] we have

$$
\begin{align*}
P \approx & I_{2}^{\diamond q_{1}} \diamond\left(\begin{array}{rr}
1 & 1 \\
0 & 1
\end{array}\right)^{\diamond q_{2}} \diamond\left(\begin{array}{rr}
1 & -1 \\
0 & 1
\end{array}\right)^{\diamond q_{3}} \diamond\left(-I_{2}\right)^{\diamond q_{4}} \\
& \diamond\left(\begin{array}{rr}
-1 & 1 \\
0 & -1
\end{array}\right)^{\diamond q_{5}} \diamond\left(\begin{array}{rr}
-1 & -1 \\
0 & -1
\end{array}\right)^{\diamond q_{6}} \\
& \diamond R\left(\theta_{1}\right) \diamond \cdots \diamond R\left(\theta_{q_{7}}\right) \diamond \cdots \diamond R\left(\theta_{q_{7}+q_{8}}\right)  \tag{5.7}\\
& \diamond N_{2}\left(\omega_{1}, b_{1}\right) \diamond \cdots \diamond N_{2}\left(\omega_{q_{9}}, b_{q_{9}}\right) \\
& \diamond D(2)^{\diamond q_{10} \diamond D(-2)^{\diamond q_{11}}} .
\end{align*}
$$

where $q_{i} \geq 0$ for $1 \leq i \leq 11$ with $q_{1}+q_{2}+\cdots+q_{8}+2 q_{9}+q_{10}+q_{11}=n$, $\theta_{j} \in(0, \pi)$ for $1 \leq j \leq q_{7}, \theta_{j} \in(\pi, 2 \pi)$ for $q_{7}+1 \leq j \leq q_{7}+q_{8}, \omega_{j} \in(\mathbf{U} \backslash \mathbb{R})$ for $1 \leq j \leq q_{9}$, and

$$
b_{j}=\left(\begin{array}{ll}
b_{j 1} & b j_{2} \\
b_{j 3} & b_{j 4}
\end{array}\right) \quad \text { satisfying } b_{j 2} \neq b_{j 3} \text { for } 1 \leq j \leq q_{9}
$$

By (5.7) and Remark 2.3 we obtain

$$
\begin{align*}
P^{2} \approx & I_{2}^{\diamond\left(q_{1}+q_{4}\right)} \diamond\left(\begin{array}{ll}
1 & 1 \\
0 & 1
\end{array}\right)^{\diamond\left(q_{2}+q_{6}\right)} \diamond\left(\begin{array}{rr}
1 & -1 \\
0 & 1
\end{array}\right)^{\diamond\left(q_{3}+q_{5}\right)} \\
& \diamond R\left(2 \theta_{1}\right) \diamond \cdots \diamond R\left(2 \theta_{q_{7}}\right) \diamond \cdots \diamond R\left(2 \theta_{q_{7}+q_{8}}\right)  \tag{5.8}\\
& \diamond N_{2}\left(\omega_{1}, b_{1}\right)^{2} \diamond \cdots \diamond N_{2}\left(\omega_{q_{9}}, b_{q_{9}}\right)^{2} \diamond D(2)^{\diamond\left(q_{10}+q_{11}\right)}
\end{align*}
$$

By theorem 7.8 of [19], (5.5), and (5.8), there hold

$$
\begin{equation*}
q_{2}+q_{6} \geq p_{1}, \quad q_{10}+q_{11} \geq p_{2} \tag{5.9}
\end{equation*}
$$

Since $\gamma^{2}(t)=\gamma(t-\tau) \gamma(\tau)$ for all $t \in[\tau, 2 \tau]$, we have $\gamma^{2}$ is also the second iteration of $\gamma$ in the periodic boundary value case, so by the Bott-type formula (cf. theorem 9.2.1 of [20]), the proof of lemma 4.1 of [21], and Lemma 2.5], we have

$$
\begin{aligned}
i\left(\gamma^{2}\right)+ & 2 S_{P^{2}}^{+}(1)-v\left(\gamma^{2}\right) \\
= & 2 i(\gamma)+2 S_{P}^{+}(1)+\sum_{\theta \in(0, \pi)}\left(S_{P}^{+}\left(e^{\sqrt{-1} \theta}\right)\right. \\
& -\left(\sum_{\theta \in(0, \pi)} S_{P}^{-}\left(e^{\sqrt{-1} \theta}\right)+\left(v(P)-S_{P}^{+}(1)\right)+\left(v_{-1}(P)-S_{P}^{-}(-1)\right)\right) \\
= & 2 i(\gamma)+2\left(q_{1}+q_{2}\right)+\left(q_{8}-q_{7}\right)-\left(q_{1}+q_{3}+q_{4}+q_{5}\right) \\
\geq & 2 n+q_{1}+2 q_{2}+\left(q_{8}-q_{7}\right)-\left(q_{3}+q_{4}+q_{5}\right) \\
= & n+\left(2 q_{1}+3 q_{2}+q_{6}+2 q_{8}+2 q_{9}+q_{10}+q_{11}\right) \\
\geq & n+2 q_{2}+q_{6}+q_{10}+q_{11} \\
\geq & n+p_{1}+p_{2}
\end{aligned}
$$

where in the first equality we have used $S_{P^{2}}^{+}(1)=S_{P}^{+}(1)+S_{P}^{+}(-1)$ and $\nu\left(\gamma^{2}\right)=$ $v(\gamma)+v_{-1}(\gamma)$, in the first inequality we have used the condition $i(\gamma) \geq n$, in the third equality we have used that $q_{1}+q_{2}+\cdots+q_{8}+2 q_{9}+q_{10}+q_{11}=n$, and in the last inequality we have used (5.9). By (5.10) Claim 5.4 holds.

We continue with the proof of Theorem 5.3. We set $\mathcal{A}=\mu_{\left(L_{1}, L_{0}\right)}(\gamma)+S_{P^{2}}^{+}(1)$ and $\mathcal{B}=\mu_{\left(L_{0}, L_{1}\right)}(\gamma)+S_{P^{2}}^{+}(1)$.

By proposition C of [21] and proposition 6.1 of [17] we have

$$
\begin{equation*}
i_{L_{0}}(\gamma)+i_{L_{1}}(\gamma)=i\left(\gamma^{2}\right)-n, \quad v_{L_{0}}(\gamma)+v_{L_{1}}(\gamma)=\nu\left(\gamma^{2}\right) \tag{5.11}
\end{equation*}
$$

From (5.11) or 5.1) we obtain

$$
\begin{equation*}
\mathcal{A}+\mathcal{B}=i\left(\gamma^{2}\right)+2 S_{P^{2}}^{+}(1)-v\left(\gamma^{2}\right)-n \tag{5.12}
\end{equation*}
$$

Case 1. $\nu_{L_{0}}(\gamma)=0$.
In this case, $i_{L_{1}}(\gamma)+S_{P^{2}}^{+}(1)-v_{L_{0}}(\gamma) \geq 0+0-0=0$ and (5.4) holds.
Case 2. $v_{L_{0}}(\gamma)=n$.
In this case

$$
P=\left(\begin{array}{ll}
A & 0 \\
C & D
\end{array}\right)
$$

so $A$ is invertible and

$$
\begin{equation*}
m^{0}\left(A^{\top} C\right)=v_{L_{1}}(P)=v_{L_{1}}(\gamma) \tag{5.13}
\end{equation*}
$$

By Lemma 4.4 we have

$$
\begin{align*}
N P^{-1} N P & =\left(\begin{array}{cc}
I_{n} & 0 \\
2 A^{\top} C & I_{n}
\end{array}\right)  \tag{5.14}\\
& \approx I_{2}^{\diamond m^{0}\left(A^{\top} C\right)} \diamond N_{1}(1,1)^{\diamond m^{-}\left(A^{\top} C\right)} \diamond N_{1}(1,-1)^{\diamond m^{+}\left(A^{\top} C\right)}
\end{align*}
$$

By Claim 5.4, (5.14), and (5.12), there holds

$$
\begin{equation*}
\mathcal{A}+\mathcal{B} \geq m^{-}\left(A^{\top} C\right) \tag{5.15}
\end{equation*}
$$

By Theorem 4.2, Lemma 4.8, and (5.13) we obtain

$$
\begin{equation*}
\mathcal{A}-\mathcal{B} \geq m^{+}\left(A^{\top} C\right)+m^{0}\left(A^{\top} C\right)-n \tag{5.16}
\end{equation*}
$$

Then 5.15) and 5.16 give

$$
2 \mathcal{A} \geq m^{-}\left(A^{\top} C\right)+\left(m^{+}\left(A^{\top} C\right)+m^{0}\left(A^{\top} C\right)\right)-n=0
$$

which yields $\mathcal{A} \geq 0$ and 5.4 holds.
Case 3. $1 \leq v_{L_{0}}(\gamma)=v_{L_{0}}(P) \leq n-1$.
In this case by (i) of Lemma 4.14 we have

$$
P:=\left(\begin{array}{cc}
A & B \\
C & D
\end{array}\right) \sim\left(\begin{array}{cccc}
A_{1} & B_{1} & I_{r} & 0 \\
0 & D_{1} & 0 & 0 \\
A_{3} & B_{3} & A_{2} & 0 \\
C_{3} & D_{3} & C_{2} & D_{2}
\end{array}\right)
$$

where $A_{1}, A_{2}, A_{3}$ are $r \times r$ matrices, $D_{1}, D_{2}, D_{3}$ are $(n-r) \times(n-r)$ matrices, $B_{1}, B_{3}$ are $r \times(n-r)$ matrices, and $C_{2}, C_{3}$ are $(n-r) \times r$ matrices. We divide Case 3 into the following three subcases.

Subcase 1. $A_{3}=0$.
In this subcase let $\lambda=\operatorname{rank} B_{3}$. Then $0 \leq \lambda \leq \min \{r, n-r\}, A_{1}$ is invertible, $A_{1} A_{2}=I_{r}$, and $D_{1} D_{2}^{\top}=I_{k-r}$, so we have $A$ is invertible; furthermore, there holds $m^{0}\left(A^{\top} C\right)=\operatorname{dim} \operatorname{ker} C=\nu_{L_{1}}(P)$. Suppose $m^{+}\left(A_{1}\right)=p, m^{-}\left(A_{1}\right)=$ $r-p$; then by (iv) of Lemma 4.14 we have

$$
\begin{align*}
& N_{k} R^{-1} N_{k} R \approx\left(\begin{array}{ll}
1 & 1 \\
0 & 1
\end{array}\right)^{\diamond p+q^{-}} \diamond\left(\begin{array}{rr}
1 & -1 \\
0 & 1
\end{array}\right)^{\diamond\left(r-p+q^{+}\right)} \diamond I_{2}^{\diamond q^{0}} \diamond D(2)^{\diamond \lambda},  \tag{5.17}\\
& m^{+}\left(A^{\top} C\right)=\lambda+q^{+},  \tag{5.18}\\
& m^{0}\left(A^{\top} C\right)=r-\lambda+q^{0},  \tag{5.19}\\
& m^{-}\left(A^{\top} C\right)=\lambda+q^{-}, \tag{5.20}
\end{align*}
$$

where $q^{*} \geq 0$ for $*=+,-, 0$ and $q^{+}+q^{0}+q^{-}=n-r-\lambda$.
By (5.17) and Claim 5.4, there holds

$$
\begin{equation*}
i\left(\gamma^{2}\right)+2 S_{P^{2}}^{+}(1)-v\left(\gamma^{2}\right) \geq n+p+q^{-}+\lambda \geq n+q^{-}+\lambda . \tag{5.21}
\end{equation*}
$$

(5.21) and (5.12) give

$$
\begin{equation*}
\mathcal{A}+\mathcal{B} \geq q^{-}+\lambda \tag{5.22}
\end{equation*}
$$

By Theorem 4.2, Lemma 4.8, and (5.18)-(5.20), we have

$$
\begin{align*}
\mathcal{A}-\mathcal{B} \geq m^{+}\left(A^{\top} C\right)+m^{0}\left(A^{\top} C\right)-n & =q^{+}+\lambda+r-\lambda+q^{0}-n \\
& =r+q^{+}+q^{0}-n . \tag{5.23}
\end{align*}
$$

Since $q^{+}+q^{0}+q^{-}=n-r-\lambda$, (5.22) and (5.23) imply

$$
\begin{aligned}
2 \mathcal{A} & \geq q^{-}+\lambda+\left(r+q^{+}+q^{0}\right)-n \\
& =\left(q^{-}+q^{+}+q^{0}\right)-(n-r-\lambda) \\
& =0,
\end{aligned}
$$

which yields (5.4).
Subcase 2. $A_{3}$ is invertible.
In this case by (ii) of Lemma 4.14 there holds

$$
P \sim\left(\begin{array}{cc}
A_{1} & I_{r}  \tag{5.24}\\
A_{3} & A_{2}
\end{array}\right) \diamond\left(\begin{array}{cc}
D_{1} & 0 \\
\tilde{D}_{3} & D_{2}
\end{array}\right):=P_{1} \diamond P_{2},
$$

where $\widetilde{D}_{3}$ is a $(k-r) \times(k-r)$ matrix. Then by (5.24) and Lemma 4.7 we obtain

$$
\begin{equation*}
P^{2} \approx\left(N_{r} P_{1}^{-1} N_{r} P_{1}\right) \diamond\left(N_{n-r} P_{2}^{-1} N_{n-r} P_{2}\right) \tag{5.25}
\end{equation*}
$$

Let $e\left(N_{r} P_{1}^{-1} N_{r} P_{1}\right)=2 m$; by Lemma 4.13 we have $0 \leq m \leq r$ and

$$
\begin{equation*}
\frac{1}{2} \operatorname{sgn} M_{\varepsilon}\left(P_{1}\right) \leq r-m, \quad 0<-\varepsilon \ll 1 \tag{5.26}
\end{equation*}
$$

Also by (5.25) and (5.8), there exists $\widetilde{P}_{1} \in \operatorname{Sp}(2 m)$ such that

$$
\begin{equation*}
N_{r} P_{1}^{-1} N_{r} P_{1} \approx D(2)^{\diamond(r-m)} \diamond \widetilde{P}_{1} \tag{5.27}
\end{equation*}
$$

By Lemma 4.4, there holds

$$
\begin{align*}
& N_{n-r} P_{2}^{-1} N_{n-r} P_{2} \\
& \quad=\left(\begin{array}{cc}
I_{n-r} & 0 \\
2 D_{1}^{\top} \widetilde{D}_{3} & I_{n-r}
\end{array}\right) \\
& \quad \approx\left(\begin{array}{ll}
1 & 1 \\
0 & 1
\end{array}\right)^{\diamond m^{-}\left(D_{1}^{\top} \widetilde{D}_{3}\right)} \diamond I_{2}^{\diamond m^{0}\left(D_{1}^{\top} \widetilde{D}_{3}\right)} \diamond\left(\begin{array}{cc}
1 & -1 \\
0 & 1
\end{array}\right)^{\diamond m^{+}\left(D_{1}^{\top} \widetilde{D}_{3}\right)} . \tag{5.28}
\end{align*}
$$

So by Claim 5.4 and (5.27), 5.28), (5.25), and 5.12) we have

$$
\begin{equation*}
\mathcal{A}+\mathcal{B} \geq m^{-}\left(D_{1}^{\top} \widetilde{D}_{3}\right)+r-m \tag{5.29}
\end{equation*}
$$

By Theorem 4.2 and Lemma 4.8 together with Lemma 4.13, for $0<-\varepsilon \ll 1$ we get

$$
\begin{align*}
\mathcal{A}-\mathcal{B} & =-\frac{1}{2} \operatorname{sgn} M_{\varepsilon}\left(P_{1}\right)-\frac{1}{2} \operatorname{sgn} M_{\varepsilon}\left(P_{2}\right) \\
& \geq-r+m-(n-r)+m^{+}\left(D_{1}^{\top} \widetilde{D}_{3}\right)+m^{0}\left(D_{1}^{\top} \widetilde{D}_{3}\right) \\
& =m+m^{+}\left(D_{1}^{\top} \widetilde{D}_{3}\right)+m^{0}\left(D_{1}^{\top} \widetilde{D}_{3}\right)-n \tag{5.30}
\end{align*}
$$

where we have used the fact that $m^{0}\left(D_{1}^{\top} \widetilde{D}_{3}\right)=\operatorname{ker}\left(\widetilde{D}_{3}\right)=v_{L_{1}}\left(P_{2}\right)$.
Note that

$$
\begin{equation*}
m^{+}\left(D_{1}^{\top} \widetilde{D}_{3}\right)+m^{0}\left(D_{1}^{\top} \widetilde{D}_{3}\right)+m^{-}\left(D_{1}^{\top} \widetilde{D}_{3}\right)=n-r \tag{5.31}
\end{equation*}
$$

Then by (5.29), 5.30), and (5.31) we have

$$
\begin{aligned}
2 \mathcal{A} & \geq m^{-}\left(D_{1}^{\top} \widetilde{D}_{3}\right)+r-m+\left(m+m^{+}\left(D_{1}^{\top} \widetilde{D}_{3}\right)+m^{0}\left(D_{1}^{\top} \widetilde{D}_{3}\right)\right)-n \\
& =m^{+}\left(D_{1}^{\top} \widetilde{D}_{3}\right)+m^{0}\left(D_{1}^{\top} \widetilde{D}_{3}\right)+m^{-}\left(D_{1}^{\top} \widetilde{D}_{3}\right)-(n-r) \\
& =0
\end{aligned}
$$

which yields (5.4).
SUBCASE 3. $1 \leq \operatorname{rank} A_{3}=l \leq r-1$.
In this case by (iii) of Lemma 4.14 there holds

$$
P \sim\left(\begin{array}{cc}
U & I_{l}  \tag{5.32}\\
\Lambda & V
\end{array}\right) \diamond\left(\begin{array}{cccc}
\tilde{A}_{1} & \widetilde{B}_{1} & I_{r-l} & 0 \\
0 & D_{1} & 0 & 0 \\
0 & \widetilde{B}_{3} & \tilde{A}_{2} & 0 \\
\widetilde{C}_{3} & \widetilde{D}_{3} & \widetilde{C}_{2} & \widetilde{D}_{2}
\end{array}\right):=P_{3} \diamond P_{4}
$$

where $\tilde{A}_{1}, \tilde{A}_{2}$ are $(r-l) \times(r-l)$ matrices, $\widetilde{B}_{1}, \widetilde{B}_{3}$ are $(r-l) \times(n-r)$ matrices, $\widetilde{C}_{2}, \widetilde{C}_{3}$ are $(n-r) \times(r-l)$ matrices, $D_{1}, \widetilde{D}_{2}, \widetilde{D}_{3}$ are $(n-r) \times(n-r)$ matrices, $U, V, \Lambda$ are $l \times l$ matrices, and $\Lambda$ is invertible.

Let $\lambda=\operatorname{rank} \widetilde{B}_{3}$ and denote

$$
P_{4}=\left(\begin{array}{cc}
\tilde{A} & \widetilde{B} \\
\widetilde{C} & \tilde{D}
\end{array}\right),
$$

where $\widetilde{A}, \widetilde{B}, \widetilde{C}, \tilde{D}$ are $(n-l)$-order real matrices. Assume $m^{+}\left(\tilde{A}_{1}\right)=p$ and $m^{-}\left(\widetilde{A}_{1}\right)=r-l-p$; then by (iv) of Lemma 4.14 we have

$$
\begin{align*}
N_{k} P_{4}^{-1} N_{k} P_{4} & \approx\left(\begin{array}{ll}
1 & 1 \\
0 & 1
\end{array}\right)^{\diamond\left(p+q^{-}\right)} \diamond\left(\begin{array}{cc}
1 & -1 \\
0 & 1
\end{array}\right)^{\diamond\left(r-l-p+q^{+}\right)}  \tag{5.33}\\
& \diamond I_{2}^{\diamond q^{0}} \diamond D(2)^{\diamond \lambda}, \\
m^{+}\left(\tilde{A}^{\top} \tilde{C}\right) & =\lambda+q^{+},  \tag{5.34}\\
m^{0}\left(\tilde{A}^{\top} \tilde{C}\right) & =r-l-\lambda+q^{0},  \tag{5.35}\\
m^{-}\left(\tilde{A}^{\top} \tilde{C}\right) & =\lambda+q^{-}, \tag{5.36}
\end{align*}
$$

where $q^{*} \geq 0$ for $*=+,-, 0$ and $q^{+}+q^{0}+q^{-}=n-r-\lambda$.
Let $e\left(N_{l} P_{3}^{-1} N_{l} P_{3}\right)=2 m$. By Lemma 4.13 we obtain $0 \leq m \leq l$ and

$$
\begin{equation*}
\frac{1}{2} \operatorname{sgn} M_{\varepsilon}\left(P_{3}\right) \leq l-m, \quad 0<-\varepsilon \ll 1 \tag{5.37}
\end{equation*}
$$

By similar argument as in the proof of Subcase 2, there exists $\widetilde{P}_{3} \in \operatorname{Sp}(2 m)$ such that

$$
\begin{equation*}
N_{r} P_{3}^{-1} N_{r} P_{3} \approx D(2)^{\diamond(l-m)} \diamond \widetilde{P}_{3} . \tag{5.38}
\end{equation*}
$$

So by Claim 5.4, (5.32), (5.33), (5.38), and (5.12) we have

$$
\begin{equation*}
\mathcal{A}+\mathcal{B} \geq q^{-}+l-m+\lambda . \tag{5.39}
\end{equation*}
$$

By Theorem 4.2, Lemma 4.8, (5.34), (5.35) and (5.37), for $0 \leq-\varepsilon \ll 1$ we obtain

$$
\begin{align*}
\mathcal{A}-\mathcal{B} & =-\frac{1}{2} \operatorname{sgn} M_{\varepsilon}\left(P_{3}\right)-\frac{1}{2} \operatorname{sgn} M_{\varepsilon}\left(P_{4}\right) \\
& \geq-\frac{1}{2} \operatorname{sgn} M_{\varepsilon}\left(P_{3}\right)-(n-l)-m^{+}\left(\tilde{A}^{\top} \widetilde{C}\right)+m^{0}\left(\tilde{A}^{\top} \widetilde{C}\right) \\
& \geq-l+m-(n-l)+\left(\lambda+q^{+}\right)+\left(r-l-\lambda+q^{0}\right) \\
& =\left(q^{+}+q^{0}+r\right)-n-(l-m) . \tag{5.40}
\end{align*}
$$

Since $q^{+}+q^{0}+q^{-}=n-r-\lambda$, by (5.39) and (5.40) we have

$$
\begin{aligned}
2 \mathcal{A} & \geq q^{-}+l-m+\lambda+\left(q^{+}+q^{0}+r\right)-n-(l-m) \\
& =\left(q^{+}+q^{0}+q^{-}\right)-(n-r-\lambda) \\
& =0,
\end{aligned}
$$

which yields (5.4). Hence (5.4) holds in Cases 1 through 3 and the proof of Theorem 5.3 is completed.

Remark 5.5. Both the estimates (5.3) and (5.4) in Theorem 5.3 are optimal . In fact, we can construct a symplectic path satisfying the conditions of Theorem 5.3 such that the equalities in (5.3) and (5.4) hold. Let $\tau=\pi$ and $\gamma(t)=R(t)^{\diamond n}, t \in$ $[0, \pi]$. It is easy to see that

$$
i_{L_{0}}(\gamma)=\sum_{0<t<\pi} v_{L_{0}}(\gamma(t))=0 \quad \text { and also } \quad i_{L_{1}}(\gamma)=\sum_{0<t<\pi} v_{L_{1}}(\gamma(t))=0,
$$

$\nu_{L_{0}}(\gamma)=\nu_{L_{1}}(\gamma)=n, \gamma^{2}(t)=\gamma(t-\pi) \gamma(\pi)$ for $t \in[\pi, 2 \pi], i(\gamma)=n$, and $P=\gamma(\pi)=-I_{2 n}$. Hence by Lemma 2.5. $S_{P^{2}}^{+}(1)=S_{I_{2 n}}^{+}(1)=n$. Thus

$$
\mu_{\left(L_{0}, L_{1}\right)}(\gamma)+S_{P^{2}}^{+}(1)=\mu_{\left(L_{1}, L_{0}\right)}(\gamma)+S_{P^{2}}^{+}(1)=0-n+n=0 .
$$

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