# Multiple brake orbits on compact convex symmetric reversible hypersurfaces in $\mathbf{R}^{2 n}$ 

Duanzhi Zhang ${ }^{1}$, Chungen Liu ${ }^{*, 2}$<br>School of Mathematics and LPMC, Nankai University, Tianjin 300071, People's Republic of China

Received 1 November 2011; accepted 26 March 2013
Available online 11 June 2013


#### Abstract

In this paper, we prove that there exist at least $\left[\frac{n+1}{2}\right]+1$ geometrically distinct brake orbits on every $C^{2}$ compact convex symmetric hypersurface $\Sigma$ in $\mathbf{R}^{2 n}$ for $n \geqslant 2$ satisfying the reversible condition $N \Sigma=\Sigma$ with $N=\operatorname{diag}\left(-I_{n}, I_{n}\right)$. As a consequence, we show that there exist at least $\left[\frac{n+1}{2}\right]+1$ geometrically distinct brake orbits in every bounded convex symmetric domain in $\mathbf{R}^{n}$ with $n \geqslant 2$ which gives a positive answer to the Seifert conjecture of 1948 in the symmetric case for $n=3$. As an application, for $n=4$ and 5 , we prove that if there are exactly $n$ geometrically distinct closed characteristics on $\Sigma$, then all of them are symmetric brake orbits after suitable time translation.


© 2013 Elsevier Masson SAS. All rights reserved.
MSC: 58E05; 70H05; 34C25
Keywords: Brake orbit; Maslov-type index; Seifert conjecture; Convex symmetric

## 1. Introduction

Let $V \in C^{2}\left(\mathbf{R}^{n}, \mathbf{R}\right)$ and $h>0$ be such that $\Omega \equiv\left\{q \in \mathbf{R}^{n} \mid V(q)<h\right\}$ is nonempty, bounded, open and connected. Consider the following fixed energy problem of the second order autonomous Hamiltonian system

$$
\begin{align*}
& \ddot{q}(t)+V^{\prime}(q(t))=0, \quad \text { for } q(t) \in \Omega,  \tag{1.1}\\
& \frac{1}{2}|\dot{q}(t)|^{2}+V(q(t))=h, \quad \forall t \in \mathbf{R},  \tag{1.2}\\
& \dot{q}(0)=\dot{q}\left(\frac{\tau}{2}\right)=0,  \tag{1.3}\\
& q\left(\frac{\tau}{2}+t\right)=q\left(\frac{\tau}{2}-t\right), \quad q(t+\tau)=q(t), \quad \forall t \in \mathbf{R} . \tag{1.4}
\end{align*}
$$

[^0]A solution $(\tau, q)$ of (1.1)-(1.4) is called a brake orbit in $\Omega$. We call two brake orbits $q_{1}$ and $q_{2}: \mathbf{R} \rightarrow \mathbf{R}^{n}$ geometrically distinct if $q_{1}(\mathbf{R}) \neq q_{2}(\mathbf{R})$.

We denote by $\mathcal{O}(\Omega)$ and $\tilde{\mathcal{O}}(\Omega)$ the sets of all brake orbits and geometrically distinct brake orbits in $\Omega$ respectively.
Let $J_{k}=\left(\begin{array}{cc}0 & -I_{k} \\ I_{k} & 0\end{array}\right)$ and $N_{k}=\left(\begin{array}{cc}-I_{k} & 0 \\ 0 & I_{k}\end{array}\right)$ with $I_{k}$ being the identity in $\mathbf{R}^{k}$. 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 $k \in \mathbf{N}$ is defined by

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

where $M^{T}$ 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

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

Suppose that $H \in C^{2}\left(\mathbf{R}^{2 n} \backslash\{0\}, \mathbf{R}\right) \cap C^{1}\left(\mathbf{R}^{2 n}, \mathbf{R}\right)$ satisfying

$$
\begin{equation*}
H(N x)=H(x), \quad \forall x \in \mathbf{R}^{2 n} \tag{1.5}
\end{equation*}
$$

We consider the following fixed energy problem

$$
\begin{align*}
& \dot{x}(t)=J H^{\prime}(x(t))  \tag{1.6}\\
& H(x(t))=h  \tag{1.7}\\
& x(-t)=N x(t)  \tag{1.8}\\
& x(\tau+t)=x(t), \quad \forall t \in \mathbf{R} \tag{1.9}
\end{align*}
$$

A solution ( $\tau, x$ ) of (1.6)-(1.9) is also called a brake orbit on $\Sigma:=\left\{y \in \mathbf{R}^{2 n} \mid H(y)=h\right\}$.

Remark 1.1. It is well known that via

$$
\begin{equation*}
H(p, q)=\frac{1}{2}|p|^{2}+V(q) \tag{1.10}
\end{equation*}
$$

$x=(p, q)$ and $p=\dot{q}$, the elements in $\mathcal{O}(\{V<h\})$ and the solutions of (1.6)-(1.9) are one-to-one correspondent.
In more general setting, let $\Sigma$ be a $C^{2}$ compact hypersurface in $\mathbf{R}^{2 n}$ bounding a compact set $C$ with nonempty interior. Suppose $\Sigma$ has non-vanishing 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 $\mathbf{R}^{2 n}$ 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 by the reversible condition there holds $N_{\Sigma}(N x)=N N_{\Sigma}(x)$. We consider the dynamics problem of finding $\tau>0$ and an absolutely continuous curve $x:[0, \tau] \rightarrow \mathbf{R}^{2 n}$ such that

$$
\begin{align*}
& \dot{x}(t)=J N_{\Sigma}(x(t)), \quad x(t) \in \Sigma,  \tag{1.11}\\
& x(-t)=N x(t), \quad x(\tau+t)=x(t), \quad \text { for all } t \in \mathbf{R} \tag{1.12}
\end{align*}
$$

A solution $(\tau, x)$ of the problem (1.11)-(1.12) is a special closed characteristic on $\Sigma$, here we still call it a brake orbit on $\Sigma$.

We also call two brake orbits $\left(\tau_{1}, x_{1}\right)$ and $\left(\tau_{2}, x_{2}\right)$ geometrically distinct if $x_{1}(\mathbf{R}) \neq x_{2}(\mathbf{R})$, otherwise we say they are equivalent. Any two equivalent brake orbits are geometrically the same. We denote by $\mathcal{J}_{b}(\Sigma)$ the set of all brake orbits on $\Sigma$, by $[(\tau, x)]$ the equivalent class of $(\tau, x) \in \mathcal{J}_{b}(\Sigma)$ in this equivalent relation and by $\tilde{\mathcal{J}}_{b}(\Sigma)$ the set of $[(\tau, x)]$ for all $(\tau, x) \in \mathcal{J}_{b}(\Sigma)$. From now on, in the notation $[(\tau, x)]$ we always assume $x$ has minimal period $\tau$. We also denote by $\tilde{\mathcal{J}}(\Sigma)$ the set of all geometrically distinct closed characteristics on $\Sigma$.

Let $(\tau, x)$ be a solution of (1.6)-(1.9). We consider the boundary value problem of the linearized Hamiltonian system

$$
\begin{align*}
& \dot{y}(t)=J H^{\prime \prime}(x(t)) y(t)  \tag{1.13}\\
& y(t+\tau)=y(t), \quad y(-t)=N y(t), \quad \forall t \in \mathbf{R} \tag{1.14}
\end{align*}
$$

Denote by $\gamma_{x}(t)$ the fundamental solution of the system (1.13), i.e., $\gamma_{x}(t)$ is the solution of the following problem

$$
\begin{align*}
& \dot{\gamma}_{x}(t)=J H^{\prime \prime}(x(t)) \gamma_{x}(t),  \tag{1.15}\\
& \gamma_{x}(0)=I_{2 n} . \tag{1.16}
\end{align*}
$$

We call $\gamma_{x} \in C([0, \tau / 2], \mathrm{Sp}(2 n))$ the associated symplectic path of $(\tau, x)$.
Let $B_{1}^{n}(0)$ denote the open unit ball $\mathbf{R}^{n}$ centered at the origin 0 . In [20] of 1948, H. Seifert proved $\tilde{\mathcal{O}}(\Omega) \neq \emptyset$ provided $V^{\prime} \neq 0$ on $\partial \Omega, V$ is analytic and $\Omega$ is homeomorphic to $B_{1}^{n}(0)$. Then he proposed his famous conjecture: $\# \tilde{O}(\Omega) \geqslant n$ under the same conditions.

After 1948, many studies have been carried out for the brake orbit problem. S. Bolotin proved first in [4] (also see [5]) of 1978 the existence of brake orbits in general setting. K. Hayashi in [10], H. Gluck and W. Ziller in [8], and V. Benci in [2] in 1983-1984 proved ${ }^{\#} \tilde{\mathcal{O}}(\Omega) \geqslant 1$ if $V$ is $C^{1}, \bar{\Omega}=\{V \leqslant h\}$ is compact, and $V^{\prime}(q) \neq 0$ for all $q \in \partial \Omega$. In 1987, P.H. Rabinowitz in [19] proved that if $H$ satisfies (1.5), $\Sigma \equiv H^{-1}(h)$ is star-shaped, and $x \cdot H^{\prime}(x) \neq 0$ for all $x \in \Sigma$, then ${ }^{\#} \tilde{\mathcal{J}}_{b}(\Sigma) \geqslant 1$. In 1987, V. Benci and F. Giannoni gave a different proof of the existence of one brake orbit in [3].

In 1989, A. Szulkin in [21] proved that ${ }^{\#} \tilde{\mathcal{J}}_{b}\left(H^{-1}(h)\right) \geqslant n$, if $H$ satisfies conditions in [19] of Rabinowitz and the energy hypersurface $H^{-1}(h)$ is $\sqrt{2}$-pinched. E.W.C. van Groesen in [9] of 1985 and A. Ambrosetti, V. Benci, Y. Long in [1] of 1993 also proved ${ }^{\#} \tilde{O}(\Omega) \geqslant n$ under different pinching conditions.

Without pinching condition, in [17] Y. Long, C. Zhu and the first author of this paper proved the following result: For $n \geqslant 2$, suppose $H$ satisfies
(H1) (smoothness) $H \in C^{2}\left(\mathbf{R}^{2 n} \backslash\{0\}, \mathbf{R}\right) \cap C^{1}\left(\mathbf{R}^{2 n}, \mathbf{R}\right)$,
(H2) (reversibility) $H(N y)=H(y)$ for all $y \in \mathbf{R}^{2 n}$,
(H3) (convexity) $H^{\prime \prime}(y)$ is positive definite for all $y \in \mathbf{R}^{2 n} \backslash\{0\}$,
(H4) (symmetry) $H(-y)=H(y)$ for all $y \in \mathbf{R}^{2 n}$.
Then for any given $h>\min \left\{H(y) \mid y \in \mathbf{R}^{2 n}\right\}$ and $\Sigma=H^{-1}(h)$, there holds

$$
{ }^{\#} \tilde{\mathcal{J}}_{b}(\Sigma) \geqslant 2 .
$$

As a consequence they also proved that: For $n \geqslant 2$, suppose $V(0)=0, V(q) \geqslant 0, V(-q)=V(q)$ and $V^{\prime \prime}(q)$ is positive definite for all $q \in \mathbf{R}^{n} \backslash\{0\}$. Then for $\Omega \equiv\left\{q \in \mathbf{R}^{n} \mid V(q)<h\right\}$ with $h>0$, there holds

$$
\# \tilde{\mathcal{O}}(\Omega) \geqslant 2
$$

Under the same condition of [17], in 2009 Liu and Zhang in [14] proved that ${ }^{\#} \tilde{\mathcal{J}}_{b}(\Sigma) \geqslant\left[\frac{n}{2}\right]+1$, also they proved \# $\tilde{\mathcal{O}}(\Omega) \geqslant\left[\frac{n}{2}\right]+1$ under the same condition of [17]. Moreover if all brake orbits on $\Sigma$ are nondegenerate, Liu and Zhang in [14] proved that ${ }^{\#} \tilde{\mathcal{J}}_{b}(\Sigma) \geqslant n+\mathfrak{A}(\Sigma)$, where $2 \mathfrak{A}(\Sigma)$ is the number of geometrically distinct asymmetric brake orbits on $\Sigma$.

Definition 1.1. We denote

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

Definition 1.2. For $\Sigma \in \mathcal{H}_{b}^{s, c}(2 n)$, a brake orbit ( $\left.\tau, x\right)$ on $\Sigma$ is called symmetric if $x(\mathbf{R})=-x(\mathbf{R})$. Similarly, for a $C^{2}$ convex symmetric bounded domain $\Omega \subset \mathbf{R}^{n}$, a brake orbit $(\tau, q) \in \mathcal{O}(\Omega)$ is called symmetric if $q(\mathbf{R})=-q(\mathbf{R})$.

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

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

The following are the main results of this paper.

Theorem 1.1. For any $\Sigma \in \mathcal{H}_{b}^{s, c}(2 n)$ with $n \geqslant 2$, we have

$$
\# \tilde{\mathcal{J}}_{b}(\Sigma) \geqslant\left[\frac{n+1}{2}\right]+1
$$

Corollary 1.1. Suppose $V(0)=0, V(q) \geqslant 0, V(-q)=V(q)$ and $V^{\prime \prime}(q)$ is positive definite for all $q \in \mathbf{R}^{n} \backslash\{0\}$ with $n \geqslant 3$. Then for any given $h>0$ and $\Omega \equiv\left\{q \in \mathbf{R}^{n} \mid V(q)<h\right\}$, we have

$$
\# \tilde{\mathcal{O}}(\Omega) \geqslant\left[\frac{n+1}{2}\right]+1
$$

Remark 1.2. Note that for $n=3$, Corollary 1.1 yields $\# \tilde{\mathcal{O}}(\Omega) \geqslant 3$, which gives a positive answer to Seifert's conjecture in the convex symmetric case.

As a consequence of Theorem 1.1, we can prove
Theorem 1.2. For $n=4,5$ and any $\Sigma \in \mathcal{H}_{b}^{s, c}(2 n)$, suppose

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

Then all of them are symmetric brake orbits after suitable translation.
Example 1.1. A typical example of $\Sigma \in \mathcal{H}_{b}^{s, c}(2 n)$ is the ellipsoid $\mathcal{E}_{n}(r)$ defined as follows. Let $r=\left(r_{1}, \ldots, r_{n}\right)$ with $r_{j}>0$ for $1 \leqslant j \leqslant n$. Define

$$
\mathcal{E}_{n}(r)=\left\{x=\left(x_{1}, \ldots, x_{n}, y_{1}, \ldots, y_{n}\right) \in \mathbf{R}^{2 n} \left\lvert\, \sum_{k=1}^{n} \frac{x_{k}^{2}+y_{k}^{2}}{r_{k}^{2}}=1\right.\right\} .
$$

If $r_{j} / r_{k} \notin \mathbf{Q}$ whenever $j \neq k$, from [7] one can see that there are precisely $n$ geometrically distinct symmetric brake orbits on $\mathcal{E}_{n}(r)$ and all of them are nondegenerate.

## 2. Index theories of $\left(i_{L_{j}}, \nu_{L_{j}}\right)$ and $\left(i_{\omega}, v_{\omega}\right)$

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

In this section, we make some preparation for the proof of Theorem 3.1 below. We first briefly review the index function $\left(i_{\omega}, v_{\omega}\right)$ and $\left(i_{L_{j}}, v_{L_{j}}\right)$ for $j=0,1$, more details can be found in [11,12,14,16]. Following Theorem 2.3 of [23] we study the differences $i_{L_{0}}(\gamma)-i_{L_{1}}(\gamma)$ and $i_{L_{0}}(\gamma)+\nu_{L_{0}}(\gamma)-i_{L_{1}}(\gamma)-\nu_{L_{1}}(\gamma)$ for $\gamma \in \mathcal{P}_{\tau}(2 n)$ by computing $\operatorname{sgn} M_{\varepsilon}(\gamma(\tau))$. We obtain some basic lemmas which will be used frequently in the proof of the main theorem of this paper.

For any $\omega \in \mathbf{U}$, the following codimension 1 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

$$
\eta * \xi(t)= \begin{cases}\xi(2 t) & \text { if } 0 \leqslant t \leqslant \frac{\tau}{2}, \\ \eta(2 t-\tau) & \text { if } \frac{\tau}{2} \leqslant t \leqslant \tau .\end{cases}
$$

Given any two $\left(2 m_{k} \times 2 m_{k}\right)$-matrices of square block form $M_{k}=\left(\begin{array}{c}A_{k} B_{k} \\ C_{k} \\ D_{k}\end{array}\right)$ for $k=1,2$, as in [16], the $\diamond$-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 $M_{1} \diamond M_{2}$ :

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

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

$$
v_{\omega}(M)=\operatorname{dim}_{\mathbf{C}} \operatorname{ker}\left(M-\omega I_{2 n}\right)
$$

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

$$
v_{\omega}(\gamma)=v_{\omega}(\gamma(\tau))
$$

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

$$
\begin{equation*}
i_{\omega}(\gamma)=\left[\operatorname{Sp}(2 n)_{\omega}^{0}: \gamma * \xi_{n}\right] \tag{2.1}
\end{equation*}
$$

where the right-hand side of (2.1) is the usual homotopy intersection number and the orientation of $\gamma * \xi_{n}$ is its positive time direction under homotopy with fixed endpoints. 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

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

Then $\left(i_{\omega}(\gamma), v_{\omega}(\gamma)\right) \in \mathbf{Z} \times\{0,1, \ldots, 2 n\}$ is called the index function of $\gamma$ at $\omega$.
For any $M \in \operatorname{Sp}(2 n)$ we define

$$
\Omega(M)=\left\{P \in \operatorname{Sp}(2 n) \mid \sigma(P) \cap \mathbf{U}=\sigma(M) \cap \mathbf{U} \text { and } \nu_{\lambda}(P)=v_{\lambda}(M), \forall \lambda \in \sigma(M) \cap \mathbf{U}\right\}
$$

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 2.1. It is easy to check that $\approx$ is an equivalent relation. If $M_{1} \approx M_{2}$, we have $M_{1}^{k} \approx M_{2}^{k}$ for any $k \in \mathbf{N}$ and $M_{1} \diamond M_{3} \approx M_{2} \diamond M_{4}$ for $M_{3} \approx M_{4}$. Also we have $P M P^{-1} \approx M$ for any $P, M \in \operatorname{Sp}(2 n)$.

The following symplectic matrices were introduced as basic normal forms in [16]:

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

where $b=\left(\begin{array}{ll}b_{1} & b_{2} \\ b_{3} & b_{4}\end{array}\right)$ with $b_{i} \in \mathbf{R}$ and $b_{2} \neq b_{3}$.
For any $M \in \operatorname{Sp}(2 n)$ and $\omega \in \mathbf{U}$, splitting number of $M$ at $\omega$ is 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 $\gamma(\tau)=M$.
Splitting numbers possesses the following properties.

Lemma 2.1. (Cf. [15], Lemma 9.1.5 and List 9.1.12 of [16].) 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
(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 $\omega \in \mathbf{U} \backslash \mathbf{R}$ 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) $\left(S_{M}^{+}(\omega), S_{M}^{-}(\omega)\right)=(1,1)$ for $\omega \in \mathbf{U} \backslash \mathbf{R}$ and $M=N_{2}(\omega, b)$ is non-trivial;
(6) $\left(S_{M}^{+}(\omega), S_{M}^{-}(\omega)\right)=(0,0)$ for any $\omega \in \mathbf{U}$ and $M \in \operatorname{Sp}(2 n)$ with $\sigma(M) \cap \mathbf{U}=\emptyset$;
(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 $\omega \in \mathbf{U}$.

Let

$$
F=\mathbf{R}^{2 n} \oplus \mathbf{R}^{2 n}
$$

possess the standard inner product. We define the symplectic structure of $F$ by

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

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

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

is a Lagrangian subspace of $F$.
Let

$$
\begin{aligned}
& V_{1}=L_{0} \times L_{0}=\{0\} \times \mathbf{R}^{n} \times\{0\} \times \mathbf{R}^{n} \subset \mathbf{R}^{4 n}, \\
& V_{2}=L_{1} \times L_{1}=\mathbf{R}^{n} \times\{0\} \times \mathbf{R}^{n} \times\{0\} \subset \mathbf{R}^{4 n} .
\end{aligned}
$$

By Proposition 6.1 of [18] and Lemma 2.8 and Definition 2.5 of [17], we give the following definition.
Definition 2.3. For any continuous path $\gamma \in \mathcal{P}_{\tau}(2 n)$, we define the following Maslov-type indices:

$$
\begin{aligned}
& i_{L_{0}}(\gamma)=\mu_{F}^{C L M}\left(V_{1}, \operatorname{Gr}(\gamma),[0, \tau]\right)-n, \\
& i_{L_{1}}(\gamma)=\mu_{F}^{C L M}\left(V_{2}, \operatorname{Gr}(\gamma),[0, \tau]\right)-n, \\
& v_{L_{j}}(\gamma)=\operatorname{dim}\left(\gamma(\tau) L_{j} \cap L_{j}\right), \quad j=0,1,
\end{aligned}
$$

where we denote by $i_{F}^{C L M}(V, W,[a, b])$ 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 $v_{L_{j}}(M)=\operatorname{dim}\left(M L_{j} \cap L_{j}\right)$.

Definition 2.4. For two paths $\gamma_{0}, \gamma_{1} \in \mathcal{P}_{\tau}(2 n)$ and $j=0,1$, we say that they are $L_{j}$-homotopic and denoted by $\gamma_{0} \sim_{L_{j}} \gamma_{1}$, if there is a continuous 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 2.2. (See [11].)
(1) If $\gamma_{0} \sim_{L_{j}} \gamma_{1}$, there hold

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

(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 there holds

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

(4) If $b_{11}(t)>0$ for any $t \in \mathbf{R}$, there holds

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

Definition 2.5. For any $\gamma \in \mathcal{P}_{\tau}$ and $k \in \mathbf{N} \equiv\{1,2, \ldots\}$, in this paper the $k$-time iteration $\gamma^{k}$ of $\gamma \in \mathcal{P}_{\tau}(2 n)$ in brake orbit boundary sense is defined by $\left.\tilde{\gamma}\right|_{[0, k \tau]}$ with

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

By [17] or Corollary 5.1 of [14] $\lim _{k \rightarrow \infty} \frac{i_{L_{0}}\left(\gamma^{k}\right)}{k}$ exists, as usual we define the mean $i_{L_{0}}$ index of $\gamma$ by $\hat{i}_{L_{0}}(\gamma)=$ $\lim _{k \rightarrow \infty} \frac{i_{L_{0}}\left(\gamma^{k}\right)}{k}$.

For any $P \in \operatorname{Sp}(2 n)$ and $\varepsilon \in \mathbf{R}$, we set

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

Then we have the following
Theorem 2.1. (See Theorem 2.3 of [23].) For $\gamma \in \mathcal{P}_{\tau}(2 k)$ with $\tau>0$, we have

$$
i_{L_{0}}(\gamma)-i_{L_{1}}(\gamma)=\frac{1}{2} \operatorname{sgn} M_{\varepsilon}(\gamma(\tau))
$$

where $\operatorname{sgn} M_{\varepsilon}(\gamma(\tau))=m^{+}\left(M_{\varepsilon}(\gamma(\tau))\right)-m^{-}\left(M_{\varepsilon}(\gamma(\tau))\right)$ is the signature of the symmetric matrix $M_{\varepsilon}(\gamma(\tau))$ and $0<\varepsilon \ll 1$. We also have

$$
\left(i_{L_{0}}(\gamma)+v_{L_{0}}(\gamma)\right)-\left(i_{L_{1}}(\gamma)+v_{L_{1}}(\gamma)\right)=\frac{1}{2} \operatorname{sign} M_{\varepsilon}(\gamma(\tau))
$$

where $0<-\varepsilon \ll 1$.
Remark 2.2. (See Remark 2.1 of [23].) For any $n_{j} \times n_{j}$ symplectic matrix $P_{j}$ with $j=1,2$ and $n_{j} \in \mathbf{N}$, 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 $\varepsilon \in \mathbf{R}$.

In the following of this section we will give some lemmas which will be used frequently in the proof of our main theorem later.

Lemma 2.3. For $k \in \mathbf{N}$ and any symplectic matrix $P=\left(\begin{array}{cc}I_{k} & 0 \\ C & I_{k}\end{array}\right)$, there holds $P \approx I_{2}^{\diamond p} \diamond N_{1}(1,1)^{\diamond q} \diamond N_{1}(1,-1)^{\diamond r}$ with $p, q, r$ satisfying

$$
m^{0}(C)=p, \quad m^{-}(C)=q, \quad m^{+}(C)=r
$$

Proof. It is clear that

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

where $B=\operatorname{diag}\left(0,-I_{m^{-}(C)}, I_{m^{+}(C)}\right)$. Since $J_{1} N_{1}(1, \pm 1)\left(J_{1}\right)^{-1}=\left(\begin{array}{cc}1 & 0 \\ \mp & 1\end{array}\right)$, by Remark 2.1 we have $N_{1}(1, \pm 1) \approx$ $\left(\begin{array}{cc}1 & 0 \\ \mp & 1\end{array}\right)$. Then

$$
P \approx I_{2}^{\diamond m^{0}(C)} \diamond N_{1}(1,1)^{\diamond m^{-}(C)} \diamond N_{1}(1,-1)^{\diamond m^{+}(C)} .
$$

By Lemma 2.1 we have

$$
\begin{equation*}
S_{P}^{+}(1)=m^{0}(C)+m^{-}(C)=p+q \tag{2.2}
\end{equation*}
$$

By the definition of the relation $\approx$, we have

$$
\begin{equation*}
2 p+q+r=v_{1}(P)=2 m^{0}(C)+m^{+}(C)+m^{-}(C) \tag{2.3}
\end{equation*}
$$

Also we have

$$
\begin{equation*}
p+q+r=m^{0}(C)+m^{+}(C)+m^{-}(C)=k \tag{2.4}
\end{equation*}
$$

By (2.2)-(2.4) we have

$$
m^{0}(C)=p, \quad m^{-}(C)=q, \quad m^{+}(C)=r
$$

The proof of Lemma 2.3 is complete.
Definition 2.6. We call two symplectic matrices $M_{1}$ and $M_{2}$ in $\operatorname{Sp}(2 k)$ special homotopic (or ( $L_{0}, L_{1}$ )-homotopic) and denote by $M_{1} \sim M_{2}$, if there are $P_{j} \in \operatorname{Sp}(2 k)$ with $P_{j}=\operatorname{diag}\left(Q_{j},\left(Q_{j}^{T}\right)^{-1}\right)$, where $Q_{j}$ is a $k \times k$ invertible real matrix, and $\operatorname{det}\left(Q_{j}\right)>0$ for $j=1,2$, such that

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

It is clear that $\sim$ is an equivalent relation.
Lemma 2.4. For $M_{1}, M_{2} \in \operatorname{Sp}(2 k)$, if $M_{1} \sim M_{2}$, then

$$
\begin{align*}
& \operatorname{sgn} M_{\varepsilon}\left(M_{1}\right)=\operatorname{sgn} M_{\varepsilon}\left(M_{2}\right), \quad 0 \leqslant|\varepsilon| \ll 1,  \tag{2.5}\\
& N_{k} M_{1}^{-1} N_{k} M_{1} \approx N_{k} M_{2}^{-1} N_{k} M_{2} . \tag{2.6}
\end{align*}
$$

Proof. By Definition 2.6, there are $P_{j} \in \operatorname{Sp}(2 k)$ with $P_{j}=\operatorname{diag}\left(Q_{j},\left(Q_{j}^{T}\right)^{-1}\right), Q_{j}$ being $k \times k$ invertible real matrix, and $\operatorname{det}\left(Q_{j}\right)>0$ such that

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

Since $\operatorname{det}\left(Q_{j}\right)>0$ for $j=1,2$, we can joint $Q_{j}$ to $I_{k}$ by invertible matrix path. Hence we can joint $P_{1} M_{2} P_{2}$ to $M_{2}$ by symplectic path preserving the nullity $\nu_{L_{0}}$ and $\nu_{L_{1}}$. By Lemma 2.2 of [23], (2.5) holds. Since $P_{j} N_{k}=N_{k} P_{j}$ for $j=1,2$. Direct computation shows that

$$
\begin{equation*}
N_{k}\left(P_{1} M_{2} P_{2}\right)^{-1} N_{k}\left(P_{1} M_{2} P_{2}\right)=P_{2}^{-1} N_{k} M_{2}^{-1} N_{k} M_{2} P_{2} \tag{2.7}
\end{equation*}
$$

Thus (2.6) holds from Remark 2.1. The proof of Lemma 2.4 is complete.

Lemma 2.5. Let $P=\left(\begin{array}{ll}A & B \\ C & D\end{array}\right) \in \operatorname{Sp}(2 k)$, where $A, B, C, D$ are all $k \times k$ matrices. Then
(i) $\frac{1}{2} \operatorname{sgn} M_{\varepsilon}(P) \leqslant k-v_{L_{0}}(P)$, for $0<\varepsilon \ll 1$. If $B=0$, we have $\frac{1}{2} \operatorname{sgn} M_{\varepsilon}(P) \leqslant 0$ for $0<\varepsilon \ll 1$.
(ii) Let $m^{+}\left(A^{T} C\right)=q$, we have

$$
\begin{equation*}
\frac{1}{2} \operatorname{sgn} M_{\varepsilon}(P) \leqslant k-q, \quad 0 \leqslant|\varepsilon| \ll 1 \tag{2.8}
\end{equation*}
$$

(iii) $\frac{1}{2} \operatorname{sgn} M_{\varepsilon}(P) \geqslant \operatorname{dim} \operatorname{ker} C-k$ for $0<\varepsilon \ll 1$. If $C=0$, then $\frac{1}{2} \operatorname{sgn} M_{\varepsilon}(P) \geqslant 0$ for $0<\varepsilon \ll 1$.
(iv) If both $B$ and $C$ are invertible, we have

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

Proof. Since $P$ is symplectic, so is for $P^{T}$. From $P^{T} J_{k} P=J_{k}$ and $P J_{k} P^{T}=J_{k}$ we get $A^{T} C, B^{T} D, A B^{T}, C D^{T}$ are all symmetric matrices and

$$
\begin{equation*}
A D^{T}-B C^{T}=I_{k}, \quad A^{T} D-C^{T} B=I_{k} \tag{2.9}
\end{equation*}
$$

We denote $s=\sin 2 \varepsilon$ and $c=\cos 2 \varepsilon$. By definition of $M_{\varepsilon}(P)$, we have

$$
\begin{array}{rl}
M_{\varepsilon}(P) & =\left(\begin{array}{cc}
A^{T} & C^{T} \\
B^{T} & D^{T}
\end{array}\right)\left(\begin{array}{cc}
s I_{k} & -c I_{k} \\
-c I_{k} & -s I_{k}
\end{array}\right)\left(\begin{array}{cc}
A & B \\
C & D
\end{array}\right)+\left(\begin{array}{cc}
s I_{k} & c I_{k} \\
c I_{k} & -s I_{k}
\end{array}\right) \\
& =\left(\begin{array}{cc}
A^{T} & C^{T} \\
B^{T} & D^{T}
\end{array}\right)\left(\begin{array}{cc}
s I_{k} & -2 c I_{k} \\
0 & -s I_{k}
\end{array}\right)\left(\begin{array}{cc}
A & B \\
C & D
\end{array}\right)+\left(\begin{array}{cc}
s I_{k} & 2 c I_{k} \\
0 & -s I_{k}
\end{array}\right) \\
& =\left(\begin{array}{cc}
s A^{T} A-2 c A^{T} C-s C^{T} C+s I_{k} & * B^{T} B-2 c B^{T} D-s D^{T} D-s I_{k}
\end{array}\right) \\
s B^{T} A-2 c B^{T} C-s D^{T} C & s B^{*} B A^{T} B-2 c C^{T} B-s C^{T} D  \tag{2.10}\\
& =\left(\begin{array}{cc}
s A^{T} A-2 c A^{T} C-s C^{T} C+s I_{k} & s B^{T} B-2 c B^{T} D-s D^{T} D-s I_{k}
\end{array}\right),
\end{array}
$$

where in the second equality we have used that $P^{T} J_{k} P=J_{k}$, in the fourth equality we have used that $M_{\varepsilon}(P)$ is a symmetric matrix. So

$$
M_{0}(P)=-2\left(\begin{array}{cc}
A^{T} C & C^{T} B \\
B^{T} C & B^{T} D
\end{array}\right)=-2\left(\begin{array}{cc}
C^{T} & 0 \\
0 & B^{T}
\end{array}\right)\left(\begin{array}{cc}
A & B \\
C & D
\end{array}\right)
$$

where we have used $A^{T} C$ is symmetric. So if both $B$ and $C$ are invertible, $M_{0}(P)$ is invertible and symmetric, its signature is invariant under small perturbation, so (iv) holds.

If $\nu_{L_{0}}(P)=\operatorname{dim} \operatorname{ker} B>0$, since $B^{T} D=D^{T} B$, for any $x \in \operatorname{ker} B \subseteq \mathbf{R}^{k}, x \neq 0$, and $0<\varepsilon \ll 1$, we have

$$
\begin{align*}
M_{\varepsilon}(P)\binom{0}{x} \cdot\binom{0}{x} & =\left(s B^{T} B-2 c D^{T} B-s D^{T} D-s I_{k}\right) x \cdot x \\
& =-s\left(D^{T} D+I_{k}\right) x \cdot x \\
& <0 \tag{2.11}
\end{align*}
$$

So $M_{\varepsilon}(P)$ is negative definite on $(0 \oplus \operatorname{ker} B) \subseteq \mathbf{R}^{2 k}$. Hence $m^{-}\left(M_{\varepsilon}(p)\right) \geqslant \operatorname{dim} \operatorname{ker} B$ which yields that $\frac{1}{2} \operatorname{sgn} M_{\varepsilon}(P) \leqslant$ $k-\operatorname{dim} \operatorname{ker} B=k-v_{L_{0}}(P)$, for $0<\varepsilon \ll 1$. Thus (i) holds. Similarly we can prove (iii).

If $m^{+}\left(A^{T} C\right)=q>0$, let $A^{T} C$ be positive definite on $E \subseteq \mathbf{R}^{k}$, then for $0 \leqslant|s| \ll 1$, similar to (2.11) we have $M_{\varepsilon}(P)$ is negative on $E \oplus 0 \subseteq \mathbf{R}^{2 k}$. Hence $m^{-}\left(M_{\varepsilon}(P)\right) \geqslant q$, which yields (2.8).

Lemma 2.6. (See [23].) For $\gamma \in \mathcal{P}_{\tau}(2), b>0$, and $0<\varepsilon \ll 1$ small enough we have
$\operatorname{sgn} M_{ \pm \varepsilon}(R(\theta))=0, \quad$ for $\theta \in \mathbf{R}$,
$\operatorname{sgn} M_{\varepsilon}(P)=0, \quad$ if $P= \pm\left(\begin{array}{cc}1 & b \\ 0 & 1\end{array}\right)$ or $\pm\left(\begin{array}{cc}1 & 0 \\ -b & 1\end{array}\right)$,

$$
\begin{aligned}
& \operatorname{sgn} M_{\varepsilon}(P)=2, \quad \text { if } P= \pm\left(\begin{array}{cc}
1 & -b \\
0 & 1
\end{array}\right) \\
& \operatorname{sgn} M_{\varepsilon}(P)=-2, \quad \text { if } P= \pm\left(\begin{array}{ll}
1 & 0 \\
b & 1
\end{array}\right) .
\end{aligned}
$$

## 3. Proofs of Theorems 1.1 and 1.2

In this section we prove Theorems 1.1 and 1.2. The proof mainly depends on the method in [14] and the following

Theorem 3.1. For any odd number $n \geqslant 3, \tau>0$ and $\gamma \in \mathcal{P}_{\tau}(2 n)$, let $P=\gamma(\tau)$. If $i_{L_{0}} \geqslant 0, i_{L_{1}} \geqslant 0$, $i(\gamma) \geqslant n$, $\gamma^{2}(t)=\gamma(t-\tau) \gamma(\tau)$ for all $t \in[\tau, 2 \tau]$, and $P \sim\left(-I_{2}\right) \diamond Q$ with $Q \in \operatorname{Sp}(2 n-2)$, then

$$
\begin{equation*}
i_{L_{1}}(\gamma)+S_{P^{2}}^{+}(1)-v_{L_{0}}(\gamma)>\frac{1-n}{2} \tag{3.1}
\end{equation*}
$$

Proof. If the conclusion of Theorem 3.1 does not hold, then

$$
\begin{equation*}
i_{L_{1}}(\gamma)+S_{P^{2}}^{+}(1)-v_{L_{0}}(\gamma) \leqslant \frac{1-n}{2} . \tag{3.2}
\end{equation*}
$$

In the following we shall obtain a contradiction from (3.2). Hence (3.1) holds and Theorem 3.1 is proved.
Since $n \geqslant 3$ and $n$ is odd, in the following of the proof of Theorem 3.1 we write $n=2 p+1$ for some $p \in \mathbf{N}$. We denote $Q=\left(\begin{array}{ll}A & B \\ C & D\end{array}\right)$, where $A, B, C, D$ are $(n-1) \times(n-1)$ matrices. Then since $Q$ is a symplectic matrix we have

$$
\begin{align*}
& A^{T} C=C^{T} A, \quad B^{T} D=D^{T} B, \quad A B^{T}=B A^{T}, \quad C D^{T}=D C^{T}  \tag{3.3}\\
& A D^{T}-B C^{T}=I_{n-1}, \quad A^{T} D-C^{T} B=I_{n-1}  \tag{3.4}\\
& \operatorname{dim} \operatorname{ker} B=v_{L_{0}}(\gamma)-1, \quad \quad \text { dimker } C=v_{L_{1}}(\gamma)-1 \tag{3.5}
\end{align*}
$$

Since $\gamma^{2}(t)=\gamma(t-\tau) \gamma(\tau)$ for all $t \in[\tau, 2 \tau]$ we have $\gamma^{2}$ is also the twice iteration of $\gamma$ in the periodic boundary value case, so by the Bott-type formula (cf. Theorem 9.2.1 of [16]) and the proof of Lemma 4.1 of [17] we have

$$
\begin{align*}
& i\left(\gamma^{2}\right)+2 S_{P^{2}}^{+}(1)-v\left(\gamma^{2}\right) \\
& \quad=2 i(\gamma)+2 S_{P}^{+}(1)+\sum_{\theta \in(0, \pi)} S_{P}^{+}\left(e^{\sqrt{-1} \theta}\right)-\sum_{\theta \in(0, \pi)} S_{P}^{-}\left(e^{\sqrt{-1} \theta}\right)+\left(v(P)-S_{P}^{-}(1)\right)+\left(v_{-1}(P)-S_{P}^{-}(-1)\right) \\
& \quad \geqslant 2 n+2 S_{P}^{+}(1)-n \\
& \quad=n+2 S_{P}^{+}(1) \\
& \quad \geqslant n \tag{3.6}
\end{align*}
$$

where we have used the condition $i(\gamma) \geqslant n$ and $S_{P^{2}}^{+}(1)=S_{P}^{+}(1)+S_{P}^{+}(-1), \nu\left(\gamma^{2}\right)=\nu(\gamma)+v_{-1}(\gamma)$. By Proposition C of [17] and Proposition 6.1 of [14] we have

$$
\begin{equation*}
i_{L_{0}}(\gamma)+i_{L_{1}}(\gamma)=i\left(\gamma^{2}\right)-n, \quad v_{L_{0}}(\gamma)+v_{L_{1}}(\gamma)=v\left(\gamma^{2}\right) \tag{3.7}
\end{equation*}
$$

So by (3.6) and (3.7) we have

$$
\begin{align*}
& \left(i_{L_{1}}(\gamma)+S_{P^{2}}^{+}(1)-v_{L_{0}}(\gamma)\right)+\left(i_{L_{0}}(\gamma)+S_{P^{2}}^{+}(1)-v_{L_{1}}(\gamma)\right) \\
& \quad=i\left(\gamma^{2}\right)+2 S_{P^{2}}^{+}(1)-v\left(\gamma^{2}\right)-n \\
& \quad \geqslant n-n \\
& \quad=0 \tag{3.8}
\end{align*}
$$

By Theorem 2.1 and Lemma 2.6 we have

$$
\begin{align*}
& \left(i_{L_{1}}(\gamma)+S_{P^{2}}^{+}(1)-v_{L_{0}}(\gamma)\right)-\left(i_{L_{0}}(\gamma)+S_{P^{2}}^{+}(1)-v_{L_{1}}(\gamma)\right) \\
& \quad=i_{L_{1}}(\gamma)-i_{L_{0}}(\gamma)-v_{L_{0}}(\gamma)+v_{L_{1}}(\gamma) \\
& \quad=-\frac{1}{2} \operatorname{sgn} M_{\varepsilon}(Q)-\frac{1}{2} \operatorname{sgn} M_{\varepsilon}\left(-I_{2}\right) \\
& \quad=-\frac{1}{2} \operatorname{sgn} M_{\varepsilon}(Q) \\
& \geqslant 1-n . \tag{3.9}
\end{align*}
$$

So by (3.8) and (3.9) we have

$$
\begin{equation*}
i_{L_{1}}(\gamma)+S_{P^{2}}^{+}(1)-v_{L_{0}}(\gamma) \geqslant \frac{1-n}{2} \tag{3.10}
\end{equation*}
$$

By (3.2), the inequality of (3.10) must be equality. Then both (3.6) and (3.9) are equality. So we have

$$
\begin{align*}
& i\left(\gamma^{2}\right)+2 S_{P^{2}}^{+}(1)-v\left(\gamma^{2}\right)=n  \tag{3.11}\\
& i_{L_{1}}(\gamma)+S_{P^{2}}^{+}(1)-v_{L_{0}}(\gamma)=\frac{1-n}{2}  \tag{3.12}\\
& i_{L_{0}}(\gamma)+v_{L_{0}}(\gamma)-i_{L_{1}}(\gamma)-v_{L_{1}}(\gamma)=n-1 . \tag{3.13}
\end{align*}
$$

Thus by (3.6), (3.11), Theorem 1.8.10 of [16], and Lemma 2.1 we have

$$
P \approx\left(-I_{2}\right)^{\diamond p_{1}} \diamond N_{1}(1,-1)^{\diamond p_{2}} \diamond N_{1}(-1,1)^{\diamond p_{3}} \diamond R\left(\theta_{1}\right) \diamond R\left(\theta_{2}\right) \diamond \cdots \diamond R\left(\theta_{p_{4}}\right),
$$

where $p_{j} \geqslant 0$ for $j=1,2,3,4, p_{1}+p_{2}+p_{3}+p_{4}=n$ and $\theta_{j} \in(0, \pi)$ for $1 \leqslant j \leqslant p_{4}$. Otherwise by (3.6) and Lemma 2.1 we have $i\left(\gamma^{2}\right)+2 S_{P^{2}}^{+}(1)-v\left(\gamma^{2}\right)>n$ which contradicts to (3.11). So by Remark 2.1, we have

$$
\begin{equation*}
P^{2} \approx I_{2}^{\diamond p_{1}} \diamond N_{1}(1,-1)^{\diamond p_{2}} \diamond R\left(\theta_{1}\right) \diamond R\left(\theta_{2}\right) \diamond \cdots \diamond R\left(\theta_{p_{3}}\right) \tag{3.14}
\end{equation*}
$$

where $p_{i} \geqslant 0$ for $1 \leqslant i \leqslant 3, p_{1}+p_{2}+p_{3}=n$ and $\theta_{j} \in(0,2 \pi)$ for $1 \leqslant j \leqslant p_{3}$.
Note that, since $\gamma^{2}(t)=\gamma(t-\tau) \gamma(\tau)$, we have

$$
\begin{equation*}
\gamma^{2}(2 \tau)=\gamma(\tau)^{2}=P^{2} . \tag{3.15}
\end{equation*}
$$

By Definition 2.5 we have

$$
\begin{equation*}
\gamma^{2}(2 \tau)=N \gamma(\tau)^{-1} N \gamma(\tau)=N P^{-1} N P . \tag{3.16}
\end{equation*}
$$

So by (3.15) and (3.16) we have

$$
\begin{equation*}
P^{2}=N P^{-1} N P \tag{3.17}
\end{equation*}
$$

By (3.17), Lemma 2.4, and $P \sim\left(-I_{2}\right) \diamond Q$ we have

$$
\begin{align*}
P^{2} & =N P^{-1} N P \\
& \approx N\left(\left(-I_{2}\right) \diamond Q\right)^{-1} N\left(\left(-I_{2}\right) \diamond Q\right) \\
& =I_{2} \diamond\left(N_{n-1} Q^{-1} N_{n-1} Q\right) . \tag{3.18}
\end{align*}
$$

So by (3.14), we have

$$
\begin{equation*}
p_{1} \geqslant 1 \tag{3.19}
\end{equation*}
$$

Also by (3.18) and Lemma 2.5, we have

$$
\begin{equation*}
P^{2} \approx I_{2} \diamond\left(N_{n-1} Q^{\prime-1} N_{n-1} Q^{\prime}\right), \quad \forall Q^{\prime} \sim Q \text { where } Q^{\prime} \in \operatorname{Sp}(2 n-2) \tag{3.20}
\end{equation*}
$$

By (3.14) it is easy to check that

$$
\begin{equation*}
\operatorname{tr}\left(P^{2}\right)=2 n-2 p_{3}+2 \sum_{j=1}^{p_{3}} \cos \theta_{j} . \tag{3.21}
\end{equation*}
$$

By (3.11), (3.14) and Lemma 2.1 we have

$$
n=i\left(\gamma^{2}\right)+2 S_{P^{2}}^{+}(1)-v\left(\gamma^{2}\right)=i\left(\gamma^{2}\right)-p_{2} \geqslant i\left(\gamma^{2}\right)-n+1 .
$$

So

$$
\begin{equation*}
i\left(\gamma^{2}\right) \leqslant 2 n-1 \tag{3.22}
\end{equation*}
$$

By (3.7) we have

$$
\begin{equation*}
i\left(\gamma^{2}\right)=n+i_{L_{0}}(\gamma)+i_{L_{1}}(\gamma) \tag{3.23}
\end{equation*}
$$

Since $i_{L_{0}}(\gamma) \geqslant 0$ and $i_{L_{1}}(\gamma) \geqslant 0$, we have $n \leqslant i\left(\gamma^{2}\right) \leqslant 2 n-1$. So we can divide the index $i\left(\gamma^{2}\right)$ into the following three cases.

Case I. $i\left(\gamma^{2}\right)=n$.
In this case, by (3.7), $i_{L_{0}}(\gamma) \geqslant 0$, and $i_{L_{1}}(\gamma) \geqslant 0$, we have

$$
\begin{equation*}
i_{L_{0}}(\gamma)=0=i_{L_{1}}(\gamma) . \tag{3.24}
\end{equation*}
$$

So by (3.13) we have

$$
\begin{equation*}
v_{L_{0}}(\gamma)-v_{L_{1}}(\gamma)=n-1 . \tag{3.25}
\end{equation*}
$$

Since $\nu_{L_{1}}(\gamma) \geqslant 1$ and $\nu_{L_{0}}(\gamma) \leqslant n$, we have

$$
\begin{equation*}
\nu_{L_{0}}(\gamma)=n, \quad \nu_{L_{1}}(\gamma)=1 . \tag{3.26}
\end{equation*}
$$

By (3.7) we have

$$
\begin{equation*}
v\left(\gamma^{2}\right)=v\left(P^{2}\right)=n+1 \tag{3.27}
\end{equation*}
$$

By (3.12), (3.24) and (3.26) we have

$$
\begin{equation*}
S_{P^{2}}^{+}(1)=\frac{1-n}{2}+n=\frac{1+n}{2}=p+1 . \tag{3.28}
\end{equation*}
$$

So by (3.14), (3.27), (3.28), and Lemma 2.1 we have

$$
\begin{equation*}
P^{2} \approx I_{2}^{\diamond(p+1)} \diamond R\left(\theta_{1}\right) \diamond \cdots \diamond R\left(\theta_{p}\right), \tag{3.29}
\end{equation*}
$$

where $\theta_{j} \in(0,2 \pi)$. By (3.5) and (3.26) we have $B=0$. By (3.18), (3.3), and (3.4), we have

$$
\begin{aligned}
P^{2} & =N P^{-1} N P \approx I_{2} \diamond\left(N_{n-1} Q^{-1} N_{n-1} Q\right) \\
& =I_{2} \diamond\left(\begin{array}{cc}
D^{T} & 0 \\
C^{T} & A^{T}
\end{array}\right)\left(\begin{array}{cc}
A & 0 \\
C & D
\end{array}\right) \\
& =I_{2} \diamond\left(\begin{array}{cc}
D^{T} A & 0 \\
2 C^{T} A & A D^{T}
\end{array}\right) \\
& =I_{2} \diamond\left(\begin{array}{cc}
I_{2 p} & 0 \\
2 A^{T} C & I_{2 p}
\end{array}\right) .
\end{aligned}
$$

Hence $\sigma\left(P^{2}\right)=\{1\}$ which contradicts to (3.29) since $p \geqslant 1$.
Case II. $i\left(\gamma^{2}\right)=n+2 k$, where $1 \leqslant k \leqslant p$.
In this case by (3.7) we have

$$
i_{L_{0}}(\gamma)+i_{L_{1}}(\gamma)=2 k
$$

Since $i_{L_{0}}(\gamma) \geqslant 0$ and $i_{L_{1}}(\gamma) \geqslant 0$ we can write $i_{L_{0}}(\gamma)=k+r$ and $i_{L_{1}}(\gamma)=k-r$ for some integer $-k \leqslant r \leqslant k$. Then by (3.13) we have

$$
\begin{equation*}
n-1 \geqslant v_{L_{0}}(\gamma)-v_{L_{1}}(\gamma)=n-2 r-1 . \tag{3.30}
\end{equation*}
$$

Thus $r \geqslant 0$ and $0 \leqslant r \leqslant k$.
By Theorem 2.1 and (i) of Lemma 2.5 we have

$$
\begin{equation*}
2 r=i_{L_{0}}(\gamma)-i_{L_{1}}(\gamma)=\frac{1}{2} M_{\varepsilon}(P) \leqslant n-v_{L_{0}}(P) \tag{3.31}
\end{equation*}
$$

which yields that $v_{L_{0}}(\gamma) \leqslant n-2 r$. So by (3.30) and $\nu_{L_{1}}(\gamma) \geqslant 1$ we have

$$
\begin{equation*}
v_{L_{0}}(\gamma)=n-2 r, \quad v_{L_{1}}(\gamma)=1 \tag{3.32}
\end{equation*}
$$

Then by (3.12) we have

$$
\begin{equation*}
S_{P^{2}}^{+}(1)=(n-2 r)+\frac{1-n}{2}-(k-r)=\frac{1+n}{2}-k-r=p+1-k-r . \tag{3.33}
\end{equation*}
$$

Then by (3.14) and $v\left(P^{2}\right)=n-2 r+1$ and Lemma 2.1 we have

$$
\begin{equation*}
P^{2} \approx I_{2}^{\diamond(p+1-k-r)} \diamond N_{1}(1,-1)^{\diamond 2 k} \diamond R\left(\theta_{1}\right) \diamond \cdots \diamond R\left(\theta_{q}\right), \tag{3.34}
\end{equation*}
$$

where $q=n-(p+1-k-r)-2 k=p+r-k \geqslant 0$. Then we have the following three subcases (i)-(iii).
(i) $q=0$.

The only possibility is $k=p$ and $r=0$, in this case $P^{2} \approx I_{2} \diamond N_{1}(1,-1)^{\diamond 2 p}$ and $B=0$. By direct computation we have

$$
N_{1}(1,-1)^{\diamond 2 p} \approx N_{2 p} Q^{-1} N_{2 p} Q=\left(\begin{array}{cc}
I_{n-1} & 0  \tag{3.35}\\
2 A^{T} C & I_{n-1}
\end{array}\right) .
$$

Then by Lemma 2.3 we have

$$
m^{+}\left(A^{T} C\right)=2 p
$$

By (ii) of Lemma 2.5 we have

$$
\begin{equation*}
\frac{1}{2} \operatorname{sgn} M_{\varepsilon}(Q) \leqslant 2 p-2 p=0, \quad 0<-\varepsilon \ll 1 . \tag{3.36}
\end{equation*}
$$

Thus by (3.36) and Theorem 2.1, for $0<-\varepsilon \ll 1$ we have

$$
\begin{aligned}
& \left(i_{L_{0}}(\gamma)+v_{L_{0}}(\gamma)\right)-\left(i_{L_{1}}(\gamma)+v_{L_{1}}(\gamma)\right) \\
& \quad=\frac{1}{2} \operatorname{sgn} M_{\varepsilon}(P) \\
& \quad=\frac{1}{2} \operatorname{sgn} M_{\varepsilon}\left(I_{2}\right)+\frac{1}{2} M_{\varepsilon}(Q) \\
& \quad=0+\frac{1}{2} M_{\varepsilon}(Q) \\
& \quad \leqslant 0
\end{aligned}
$$

which contradicts (3.13).
(ii) $q>0$ and $r=0$.

In this case $\nu_{L_{0}}(\gamma)=n$ and $\nu_{L_{1}}(\gamma)=1$, also we have $B=0$. By the equality of (3.35) we have

$$
\operatorname{tr}\left(P^{2}\right)=2 n
$$

which contradicts to (3.21) with $p_{3}=q>0$.
(iii) $q>0$ and $r>0$.

In this case, by (3.33) we have $r<p$. (Otherwise, then $p=r=k$. From (3.19) there holds $S_{P^{2}}^{+}(1) \geqslant 1$, so from (3.33) we have $1 \leqslant S_{P^{2}}^{+}(1)=1-p \leqslant 0$ a contradiction.) Here it is easy to see rank $B=2 r$. Then there are two invertible $2 p \times 2 p$ matrices $U$ and $V$ with $\operatorname{det} U>0$ and $\operatorname{det} V>0$ such that

$$
U B V=\left(\begin{array}{cc}
I_{2 r} & 0 \\
0 & 0
\end{array}\right) .
$$

So there holds

$$
Q \sim \operatorname{diag}\left(U,\left(U^{T}\right)^{-1}\right) Q \operatorname{diag}\left(\left(V^{T}\right)^{-1}, V\right)=\left(\begin{array}{cccc}
A_{1} & B_{1} & I_{2 r} & 0  \tag{3.37}\\
C_{1} & D_{1} & 0 & 0 \\
A_{3} & B_{3} & A_{2} & B_{2} \\
C_{3} & D_{3} & C_{2} & D_{2}
\end{array}\right):=Q_{1},
$$

where for $j=1,2,3, A_{j}$ is a $2 r \times 2 r$ matrix, $D_{j}$ is a $(2 p-2 r) \times(2 p-2 r)$ matrix for $j=1,2,3, B_{j}$ is a $2 r \times(2 p-2 r)$ matrix, and $C_{j}$ is $(2 p-2 r) \times 2 r$ matrix. Since $Q_{1}$ is still a symplectic matrix, we have $Q_{1}^{T} J_{2 p} Q_{1}=J_{2 p}$, then it is easy to check that

$$
\begin{equation*}
C_{1}=0, \quad B_{2}=0 \tag{3.38}
\end{equation*}
$$

So

$$
Q_{1}=\left(\begin{array}{cccc}
A_{1} & B_{1} & I_{2 r} & 0  \tag{3.39}\\
0 & D_{1} & 0 & 0 \\
A_{3} & B_{3} & A_{2} & 0 \\
C_{3} & D_{3} & C_{2} & D_{2}
\end{array}\right) .
$$

So for the case (iii) of Case II, we have the following Subcases 1-3.
Subcase 1. $A_{3}=0$.
In this case since $Q_{1}$ is symplectic, by direct computation we have

$$
N_{2 p} Q_{1}^{-1} N_{2 p} Q_{1}=\left(\begin{array}{cccc}
I_{2 r} & * & * & * \\
* & I_{2 p-2 r} & * & * \\
* & * & I_{2 r} & * \\
* & * & * & I_{2 p-2 r}
\end{array}\right)
$$

Hence we have

$$
\operatorname{tr}\left(N_{2 p} Q_{1}^{-1} N_{2 p} Q_{1}\right)=4 p
$$

Since $Q_{1} \sim Q$, we have

$$
\begin{equation*}
P \sim\left(-I_{2}\right) \diamond Q_{1} . \tag{3.40}
\end{equation*}
$$

Then by the proof of Lemma 2.4 we have

$$
\begin{align*}
\operatorname{tr} P^{2} & =\operatorname{tr}\left(N P^{-1} N P\right) \\
& =\operatorname{tr} N\left(\left(-I_{2}\right) \diamond Q_{1}\right)^{-1} N\left(\left(-I_{2}\right) \diamond Q_{1}\right) \\
& =\operatorname{tr} I_{2} \diamond\left(N_{2 p} Q_{1}^{-1} N_{2 p} Q_{1}\right) \\
& =4 p+2=2 n . \tag{3.41}
\end{align*}
$$

By (3.21) and $p_{3}=q>0$ we have

$$
\begin{equation*}
\operatorname{tr}\left(P^{2}\right)<2 n . \tag{3.42}
\end{equation*}
$$

(3.41) and (3.42) yield a contradiction.

Subcase 2. $A_{3}$ is invertible.
By $Q_{1}$ is symplectic we have

$$
\left(\begin{array}{cc}
A_{1}^{T} & 0  \tag{3.43}\\
B_{1}^{T} & D_{1}^{T}
\end{array}\right)\left(\begin{array}{cc}
A_{2} & 0 \\
C_{2} & D_{2}
\end{array}\right)-\left(\begin{array}{cc}
A_{3}^{T} & C_{3}^{T} \\
B_{3}^{T} & D_{3}^{T}
\end{array}\right)\left(\begin{array}{cc}
I_{2 r} & 0 \\
0 & 0
\end{array}\right)=I_{2 p} .
$$

Hence

$$
\begin{equation*}
D_{1}^{T} D_{2}=I_{2 p-2 r} \tag{3.44}
\end{equation*}
$$

By direct computation we have

$$
\left(\begin{array}{cccc}
A_{1} & B_{1} & I_{2 r} & 0 \\
0 & D_{1} & 0 & 0 \\
A_{3} & B_{3} & A_{2} & 0 \\
C_{3} & D_{3} & C_{2} & D_{2}
\end{array}\right)\left(\begin{array}{cccc}
I_{2 r} & -A_{3}^{-1} B_{3} & 0 & 0 \\
0 & I_{2 p-2 r} & 0 & 0 \\
0 & 0 & I_{2 r} & 0 \\
0 & 0 & B_{3}^{T}\left(A_{3}^{T}\right)^{-1} & I_{2 p-2 r}
\end{array}\right)=\left(\begin{array}{cccc}
A_{1} & \tilde{B}_{1} & I_{2 r} & 0 \\
0 & D_{1} & 0 & 0 \\
A_{3} & 0 & A_{2} & 0 \\
C_{3} & \tilde{D}_{3} & \tilde{C}_{2} & D_{2}
\end{array}\right) .
$$

So by (3.44) we have

$$
\left(\begin{array}{cccc}
I_{2 r} & -\tilde{B}_{1} D_{2}^{T} & 0 & 0 \\
0 & I_{2 p-2 r} & 0 & 0 \\
0 & 0 & I_{2 r} & 0 \\
0 & 0 & D_{2} \tilde{B}_{1}^{T} & I_{2 p-2 r}
\end{array}\right)\left(\begin{array}{cccc}
A_{1} & \tilde{B}_{1} & I_{2 r} & 0 \\
0 & D_{1} & 0 & 0 \\
A_{3} & 0 & A_{2} & 0 \\
C_{3} & \tilde{D}_{3} & \tilde{C}_{2} & D_{2}
\end{array}\right)=\left(\begin{array}{cccc}
A_{1} & 0 & I_{2 r} & 0 \\
0 & D_{1} & 0 & 0 \\
A_{3} & 0 & A_{2} & 0 \\
\tilde{C}_{3} & \tilde{D}_{3} & \hat{C}_{2} & D_{2}
\end{array}\right):=Q_{2}
$$

Then we have

$$
\begin{equation*}
Q_{2} \sim Q_{1} \sim Q \tag{3.45}
\end{equation*}
$$

Since $Q_{2}$ is a symplectic matrix, we have $Q_{2}^{T} J_{2 p} Q_{2}=J_{2 p}$, then it is easy to check that

$$
\begin{equation*}
\tilde{C}_{3}=0, \quad \hat{C}_{2}=0 \tag{3.46}
\end{equation*}
$$

Hence we have

$$
Q_{2}=\left(\begin{array}{cc}
A_{1} & I_{2 r}  \tag{3.47}\\
A_{3} & A_{2}
\end{array}\right) \diamond\left(\begin{array}{cc}
D_{1} & 0 \\
\tilde{D}_{3} & D_{2}
\end{array}\right) .
$$

Since

$$
N_{2 p-2 r}\left(\begin{array}{cc}
D_{1} & 0  \tag{3.48}\\
\tilde{D}_{3} & D_{2}
\end{array}\right)^{-1} N_{2 p-2 r}\left(\begin{array}{cc}
D_{1} & 0 \\
\tilde{D}_{3} & D_{2}
\end{array}\right)=\left(\begin{array}{cc}
I_{2 p-2 r} & 0 \\
2 D_{1}^{T} \tilde{D}_{3} & I_{2 p-2 r}
\end{array}\right),
$$

by (3.45), (3.20), and Lemma 2.4, there is a symplectic matrix $W$ such that

$$
P^{2} \approx I_{2} \diamond W \diamond\left(\begin{array}{cc}
I_{2 p-2 r} & 0  \tag{3.49}\\
2 D_{1}^{T} \tilde{D}_{3} & I_{2 p-2 r}
\end{array}\right) .
$$

Then by (3.14) and Lemma 2.3, $D_{1}^{T} \tilde{D}_{3}$ is semipositive and

$$
1+m^{0}\left(D_{1}^{T} \tilde{D}_{3}\right) \leqslant S_{P^{2}}^{+}(1)
$$

So by (3.33) we have

$$
\begin{equation*}
m^{0}\left(D_{1}^{T} \tilde{D}_{3}\right) \leqslant p+1-k-r-1=p-k-r=(2 p-2 r)-(p+k-r) \leqslant 2 p-2 r-1 . \tag{3.50}
\end{equation*}
$$

Since $D_{1}^{T} \tilde{D}_{3}$ is a semipositive $(2 p-2 r) \times(2 p-2 r)$ matrix, by (3.50) we have $m^{+}\left(D_{1}^{T} \tilde{D}_{3}\right)>0$. Then by Theorem 2.1,
(ii) of Lemma 2.5 and Lemma 2.6, for $0<-\varepsilon \ll 1$ we have

$$
\begin{align*}
& \left(i_{L_{0}}(\gamma)+v_{L_{0}}(\gamma)\right)-\left(i_{L_{1}}(\gamma)+v_{L_{1}}(\gamma)\right) \\
& \quad=\frac{1}{2}\left(M_{\varepsilon}\left(-I_{2}\right)+M_{\varepsilon}\left(\left(\begin{array}{cc}
A_{1} & I_{2 r} \\
A_{3} & A_{2}
\end{array}\right)\right)+M_{\varepsilon}\left(\left(\begin{array}{cc}
D_{1} & 0 \\
\tilde{D}_{3} & D_{2}
\end{array}\right)\right)\right) \\
& \quad \leqslant \frac{1}{2}(0+4 r+2(2 p-2 r-1)) \\
& \quad=2 p-1 \\
& \quad=n-2 \tag{3.51}
\end{align*}
$$

which contradicts to (3.13).
Subcase 3. $A_{3} \neq 0$ and $A_{3}$ is not invertible.
In this case, suppose rank $A_{3}=\lambda$, then $0<\lambda<2 r$. There is an invertible $2 r \times 2 r$ matrix $G$ with $\operatorname{det} G>0$ such that

$$
G A_{3} G^{-1}=\left(\begin{array}{cc}
\Lambda & 0  \tag{3.52}\\
0 & 0
\end{array}\right)
$$

where $\Lambda$ is a $\lambda \times \lambda$ invertible matrix. Then we have

$$
\left.\left.\begin{array}{l}
\left(\begin{array}{cccc}
\left(G^{T}\right)^{-1} & 0 & 0 & 0 \\
0 & I_{2 p-2 r} & 0 & 0 \\
0 & 0 & G & 0 \\
0 & 0 & 0 & I_{2 p-2 r}
\end{array}\right)\left(\begin{array}{cccc}
A_{1} & B_{1} & I_{2 r} & 0 \\
0 & D_{1} & 0 & 0 \\
A_{3} & B_{3} & A_{2} & 0 \\
C_{3} & D_{3} & C_{2} & D_{2}
\end{array}\right)\left(\begin{array}{ccc}
(G)^{-1} & 0 & 0 \\
0 & I_{2 p-2 r} & 0 \\
0 & 0 & G^{T}
\end{array} 0^{0}\right. \\
0  \tag{3.53}\\
0
\end{array}\right) 0 \begin{array}{l}
I_{2 p-2 r}
\end{array}\right) .
$$

By (3.52) we can write $Q_{3}$ as the following block form

$$
Q_{3}=\left(\begin{array}{cccccc}
U_{1} & U_{2} & F_{1} & I_{\lambda} & 0 & 0  \tag{3.54}\\
U_{3} & U_{4} & F_{2} & 0 & I_{2 r-\lambda} & 0 \\
0 & 0 & D_{1} & 0 & 0 & 0 \\
\Lambda & 0 & E_{1} & W_{1} & W_{2} & 0 \\
0 & 0 & E_{2} & W_{3} & W_{4} & 0 \\
G_{1} & G_{2} & D_{3} & K_{1} & K_{2} & D_{2}
\end{array}\right)
$$

Let $R_{1}=\left(\begin{array}{ccc}I_{\lambda} & 0 & 0 \\ 0 & I_{2 r-\lambda} & 0 \\ -G_{1} \Lambda^{-1} & 0 & I_{2 p-2 r}\end{array}\right)$ and $R_{2}=\left(\begin{array}{ccc}I_{\lambda} & 0 & -\Lambda^{-1} E_{1} \\ 0 & I_{2 r-\lambda} & 0 \\ 0 & 0 & I_{2 p-2 r}\end{array}\right)$. By (3.54) we have

$$
\operatorname{diag}\left(\left(R_{1}^{T}\right)^{-1}, R_{1}\right) Q_{3} \operatorname{diag}\left(R_{2},\left(R_{2}^{T}\right)^{-1}\right)=\left(\begin{array}{cccccc}
U_{1} & U_{2} & \tilde{F}_{1} & I_{\lambda} & 0 & 0 \\
U_{3} & U_{4} & \tilde{F}_{2} & 0 & I_{2 r-\lambda} & 0 \\
0 & 0 & D_{1} & 0 & 0 & 0 \\
\Lambda & 0 & 0 & W_{1} & W_{2} & 0 \\
0 & 0 & E_{2} & W_{3} & W_{4} & 0 \\
0 & G_{2} & \tilde{D}_{3} & \tilde{K}_{1} & \tilde{K}_{2} & D_{2}
\end{array}\right):=Q_{4}
$$

Since $Q_{4}$ is a symplectic matrix we have

$$
Q_{4}^{T} J Q_{4}=J
$$

Then by (3.55) and direct computation we have $U_{2}=0, U_{3}=0, W_{2}=0, W_{3}=0, \tilde{F}_{1}=0, \tilde{K}_{1}=0$, and $U_{1}, U_{4}, W_{1}$, $W_{4}$ are all symmetric matrices, and

$$
\begin{align*}
& U_{4} W_{4}=I_{2 r-\lambda},  \tag{3.55}\\
& D_{1} D_{2}^{T}=I_{2 p-2 r},  \tag{3.56}\\
& U_{4} \tilde{E}_{2}=G_{2}^{T} D_{1}, \tag{3.57}
\end{align*}
$$

So

$$
Q_{4}=\left(\begin{array}{cccccc}
U_{1} & 0 & 0 & I_{\lambda} & 0 & 0  \tag{3.58}\\
0 & U_{4} & \tilde{F}_{2} & 0 & I_{2 r-\lambda} & 0 \\
0 & 0 & D_{1} & 0 & 0 & 0 \\
\Lambda & 0 & 0 & W_{1} & 0 & 0 \\
0 & 0 & \tilde{E}_{2} & 0 & W_{4} & 0 \\
0 & G_{2} & \tilde{D}_{3} & 0 & K_{2} & D_{2}
\end{array}\right) .
$$

By (3.55)-(3.57), we have both $\tilde{E}_{2}$ and $G_{2}$ are zero or nonzero. By Definition 2.3 we have $Q_{4} \sim Q_{3} \sim Q$. Then by (3.32), $\left(\begin{array}{ccc}\Lambda & 0 & 0 \\ 0 & 0 & \tilde{E}_{2} \\ 0 & G_{2} & \tilde{D}_{3}\end{array}\right)$ is invertible. So both $\tilde{E}_{2}$ and $G_{2}$ are nonzero.

Since $Q_{4}$ is symplectic, by (3.57) we have

$$
\left(\begin{array}{ccc}
U_{1} & 0 & 0  \tag{3.59}\\
0 & U_{4} & \tilde{F}_{2} \\
0 & 0 & D_{1}
\end{array}\right)^{T}\left(\begin{array}{ccc}
\Lambda & 0 & 0 \\
0 & 0 & \tilde{E}_{2} \\
0 & G_{2} & \tilde{D}_{3}
\end{array}\right)=\left(\begin{array}{ccc}
U_{1} \Lambda & 0 & 0 \\
0 & 0 & U_{4} \tilde{E}_{2} \\
0 & \left(U_{4} \tilde{E}_{2}\right)^{T} & D_{1}^{T} \tilde{D}_{3}+\tilde{B}_{2}^{T} \tilde{E}_{2}
\end{array}\right)
$$

which is a symmetric matrix.
Denote $F=\left(\begin{array}{cc}0 & U_{4} \tilde{E}_{2} \\ \left(U_{4} \tilde{E}_{2}\right)^{T} & D_{1}^{T} \tilde{D}_{3}+\tilde{B}_{2}^{T} \tilde{E}_{2}\end{array}\right)$. Since $U_{4} \tilde{E}_{2}$ is nonzero, in the following we prove that $m^{+}(F) \geqslant 1$.
Note that here $U_{4} \tilde{E}_{2}$ is a $(2 r-\lambda) \times(2 p-2 r)$ matrix and $D_{1}^{T} \tilde{D}_{3}+\tilde{B}_{2}^{T} \tilde{E}_{2}$ is a $(2 p-2 r) \times(2 p-2 r)$ matrix. Denote $U_{4} \tilde{E}_{2}=\left(e_{i j}\right)$ and $D_{1}^{T} \tilde{D}_{3}+\tilde{B}_{2}^{T} \tilde{E}_{2}=\left(d_{i j}\right)$, where $e_{i j}$ and $d_{i j}$ are elements on the $i$-th row and $j$-th column of the corresponding matrix. Since $U_{4} \tilde{E}_{2}$ is nonzero, there exists an $e_{i j} \neq 0$ for some $1 \leqslant i \leqslant 2 r-\lambda$ and $1 \leqslant j \leqslant 2 p-2 r$. Let $x=\left(0, \ldots, 0, e_{i j}, 0, \ldots, 0\right)^{T} \in \mathbf{R}^{2 r-\lambda}$ whose $i$-th row is $e_{i j}$ and other rows are all zero, and $y=(0, \ldots, 0, \rho, 0, \ldots, 0)^{T} \in \mathbf{R}^{2 p-2 r}$ whose $j$-th row is $\rho$ and other rows are all zero. Then we have

$$
F\binom{x}{y} \cdot\binom{x}{y}=2 \rho e_{i j}^{2}-\rho^{2} d_{j j}>0
$$

for $\rho>0$ small enough. Hence the dimension of positive definite space of $F$ is at least 1 , thus $m^{+}(F) \geqslant 1$. Then

$$
m^{+}\left(\left(\begin{array}{ccc}
U_{1} \Lambda & 0 & 0  \tag{3.60}\\
0 & 0 & U_{4} \tilde{E}_{2} \\
0 & \left(U_{4} \tilde{E}_{2}\right)^{T} & D_{1}^{T} \tilde{D}_{3}+\tilde{B}_{2}^{T} \tilde{E}_{2}
\end{array}\right)\right)=m^{+}(\Lambda)+m^{+}(F) \geqslant 1 .
$$

Then by (3.59), (3.60) and (ii) of Lemma 2.5, we have

$$
\begin{equation*}
\frac{1}{2} \operatorname{sgn} M_{\varepsilon}\left(Q_{4}\right) \leqslant 2 p-1=n-2, \quad 0<-\varepsilon \ll 1 . \tag{3.61}
\end{equation*}
$$

Since $Q \sim Q_{4}$, by (3.61) and Lemma 2.4 we have

$$
\begin{equation*}
\frac{1}{2} \operatorname{sgn} M_{\varepsilon}(Q) \leqslant 2 p-1, \quad 0<-\varepsilon \ll 1 . \tag{3.62}
\end{equation*}
$$

Then since $P \sim\left(-I_{2}\right) \diamond Q$, by Theorem 2.1, Remark 2.2 and Lemma 2.4 we have

$$
\begin{aligned}
& \left(i_{L_{0}}(\gamma)+v_{L_{0}}(\gamma)\right)-\left(i_{L_{1}}(\gamma)+v_{L_{1}}(\gamma)\right) \\
& \quad=\frac{1}{2} M_{\varepsilon}(P) \\
& \quad=\frac{1}{2} \operatorname{sgn} M_{\varepsilon}\left(\left(-I_{2}\right) \diamond Q\right) \\
& \quad=\frac{1}{2} \operatorname{sgn} M_{\varepsilon}\left(-I_{2}\right)+\frac{1}{2} \operatorname{sgn} M_{\varepsilon}(Q)
\end{aligned}
$$

$$
\begin{align*}
& =0+\frac{1}{2} \operatorname{sgn} M_{\varepsilon}(Q) \\
& \leqslant n-2 \tag{3.63}
\end{align*}
$$

Thus (3.13) and (3.63) yields a contradiction. And in Case II we can always obtain a contradiction.
Case III. $i\left(\gamma^{2}\right)=n+2 k+1$, where $0 \leqslant k \leqslant p-1$.
In this case by (3.7) we have

$$
\begin{equation*}
i_{L_{0}}(\gamma)+i_{L_{1}}(\gamma)=2 k+1 \tag{3.64}
\end{equation*}
$$

Since $i_{L_{0}}(\gamma) \geqslant 0$ and $i_{L_{1}}(\gamma) \geqslant 0$ we can write $i_{L_{0}}(\gamma)=k+1+r$ and $i_{L_{1}}(\gamma)=k-r$ for some integer $-k \leqslant r \leqslant k$. Then by (3.13) we have

$$
\begin{equation*}
n-1 \geqslant v_{L_{0}}(\gamma)-v_{L_{1}}(\gamma)=n-2 r-2 \tag{3.65}
\end{equation*}
$$

Thus $r \geqslant 0$ and $0 \leqslant r \leqslant k$.
By Theorem 2.1 and (i) of Lemma 2.5 we have

$$
\begin{equation*}
2 r+1=i_{L_{0}}(\gamma)-i_{L_{1}}(\gamma)=\frac{1}{2} M_{\varepsilon}(P) \leqslant n-v_{L_{0}}(\gamma) \tag{3.66}
\end{equation*}
$$

which yields $\nu_{L_{0}}(\gamma) \leqslant n-2 r-1$. Then by (3.65) and $\nu_{L_{1}}(\gamma) \geqslant 1$ we have

$$
\begin{equation*}
v_{L_{0}}(\gamma)=n-2 r-1, \quad v_{L_{1}}(\gamma)=1 \tag{3.67}
\end{equation*}
$$

Then by (3.12) we have

$$
\begin{equation*}
S_{P^{2}}^{+}(1)=(n-2 r-1)+\frac{1-n}{2}-(k-r)=\frac{1+n}{2}-k-r-1=p-k-r \geqslant 1 \tag{3.68}
\end{equation*}
$$

Then by (3.14) and $\nu\left(P^{2}\right)=v_{L_{0}}(\gamma)+v_{L_{1}}(\gamma)=n-2 r$ and Lemma 2.1 we have

$$
P^{2} \approx I_{2}^{\diamond(p-k-r)} \diamond N_{1}(1,-1)^{\diamond(2 k+1)} \diamond R\left(\theta_{1}\right) \diamond \cdots \diamond R\left(\theta_{q}\right)
$$

where $q=n-(p-k-r)-(2 k+1)=p+r-k \geqslant p-k \geqslant 1$.
Since in this case rank $B=2 r+1 \leqslant n-2$, by the same argument of (iii) in Case II, we have

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

Then by the same argument of Subcases $1,2,3$ of Case II, we can always obtain a contradiction in Case III. The proof of Theorem 3.1 is complete.

Now we are ready to give a proof of Theorem 1.1. 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

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

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

$$
\begin{equation*}
H_{\alpha}(x)=\left(j_{\Sigma}(x)\right)^{\alpha}, \quad \alpha>1, \quad H_{\Sigma}(x)=H_{2}(x), \quad \forall x \in \mathbf{R}^{2 n} \tag{3.69}
\end{equation*}
$$

Then $H_{\Sigma} \in C^{2}\left(\mathbf{R}^{2 n} \backslash\{0\}, \mathbf{R}\right) \cap C^{1,1}\left(\mathbf{R}^{2 n}, \mathbf{R}\right)$.

We consider the following fixed energy problem

$$
\begin{align*}
& \dot{x}(t)=J H_{\Sigma}^{\prime}(x(t))  \tag{3.70}\\
& H_{\Sigma}(x(t))=1  \tag{3.71}\\
& x(-t)=N x(t),  \tag{3.72}\\
& x(\tau+t)=x(t), \quad \forall t \in \mathbf{R} . \tag{3.73}
\end{align*}
$$

Denote by $\mathcal{J}_{b}(\Sigma, 2)\left(\mathcal{J}_{b}(\Sigma, \alpha)\right.$ for $\alpha=2$ in (3.69)) the set of all solutions $(\tau, x)$ of problem (3.70)-(3.73) and by $\tilde{\mathcal{J}}_{b}(\Sigma, 2)$ the set of all geometrically distinct solutions of (3.70)-(3.73). By Remark 1.2 of [14] or discussion in [17], elements in $\mathcal{J}_{b}(\Sigma)$ and $\mathcal{J}_{b}(\Sigma, 2)$ are one-to-one correspondent. So we have ${ }^{\#} \tilde{\mathcal{J}}_{b}(\Sigma)={ }^{\#} \tilde{\mathcal{J}}_{b}(\Sigma, 2)$.

For readers' convenience in the following we list some known results which will be used in the proof of Theorem 1.1. In the following of this paper, 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 $k \in \mathbf{N}$, where $\gamma^{k}$ is defined by Definition 2.5. We have

Lemma 3.1. (See Theorem 1.5 of [14] and Theorem 4.3 of [18].) Let $\gamma_{j} \in \mathcal{P}_{\tau_{j}}(2 n)$ for $j=1, \ldots, q$. Let $M_{j}=$ $\gamma_{j}^{2}\left(2 \tau_{j}\right)=N \gamma_{j}\left(\tau_{j}\right)^{-1} N \gamma_{j}\left(\tau_{j}\right)$, 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 $\left(R, m_{1}, m_{2}, \ldots, m_{q}\right) \in \mathbf{N}^{q+1}$ such that
(i) $\nu_{L_{0}}\left(\gamma_{j}, 2 m_{j} \pm 1\right)=v_{L_{0}}\left(\gamma_{j}\right)$,
(ii) $i_{L_{0}}\left(\gamma_{j}, 2 m_{j}-1\right)+v_{L_{0}}\left(\gamma_{j}, 2 m_{j}-1\right)=R-\left(i_{L_{1}}\left(\gamma_{j}\right)+n+S_{M_{j}}^{+}(1)-v_{L_{0}}\left(\gamma_{j}\right)\right)$,
(iii) $i_{L_{0}}\left(\gamma_{j}, 2 m_{j}+1\right)=R+i_{L_{0}}\left(\gamma_{j}\right)$.
and
(iv) $\nu\left(\gamma_{j}^{2}, 2 m_{j} \pm 1\right)=v\left(\gamma_{j}^{2}\right)$,
(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 $i\left(\gamma_{j}^{2}, n_{j}\right)=i\left(\gamma_{j}^{2 n_{j}},\left[0,2 n_{j} \tau_{j}\right]\right), \nu\left(\gamma_{j}^{2}, n_{j}\right)=\nu\left(\gamma_{j}^{2 n_{j}},\left[0,2 n_{j} \tau_{j}\right]\right)$ for $n_{j} \in \mathbf{N}$.
Lemma 3.2. (See Lemma 1.1 of [14].) Let $(\tau, x) \in \mathcal{J}_{b}(\Sigma, 2)$ be symmetric in the sense that $x\left(t+\frac{\tau}{2}\right)=-x(t)$ for all $t \in \mathbf{R}$ and $\gamma$ be the associated symplectic path of $(\tau, x)$. Set $M=\gamma\left(\frac{\tau}{2}\right)$. Then there is a continuous symplectic path

$$
\Psi(s)=P(s) M P(s)^{-1}, \quad s \in[0,1],
$$

such that

$$
\begin{aligned}
& \Psi(0)=M, \quad \Psi(1)=\left(-I_{2}\right) \diamond \tilde{M}, \quad \tilde{M} \in S p(2 n-2), \\
& v_{1}(\Psi(s))=v_{1}(M), \quad v_{2}(\Psi(s))=v_{2}(M), \quad \forall s \in[0,1],
\end{aligned}
$$

where $P(s)=\left(\begin{array}{cc}\psi(s)^{-1} & 0 \\ 0 & \psi(s)^{T}\end{array}\right)$ and $\psi$ is a continuous $n \times n$ matrix path with $\operatorname{det} \psi(s)>0$ for all $s \in[0,1]$.
For any $(\tau, x) \in \mathcal{J}_{b}(\Sigma, 2)$ and $m \in \mathbf{N}$, as in [14] we denote $i_{L_{j}}(x, m)=i_{L_{j}}\left(\gamma_{x}^{m},\left[0, \frac{m \tau}{2}\right]\right)$ and $\nu_{L_{j}}(x, m)=$ $\nu_{L_{j}}\left(\gamma_{x}^{m},\left[0, \frac{m \tau}{2}\right]\right)$ for $j=0,1$ respectively. Also we denote $i(x, m)=i\left(\gamma_{x}^{2 m},[0, m \tau]\right)$ and $\nu(x, m)=v\left(\gamma_{x}^{2 m},[0, m \tau]\right)$. If $m=1$, we denote $i(x)=i(x, 1)$ and $\nu(x)=v(x, 1)$. By Lemma 6.3 of [14] we have

Lemma 3.3. Suppose ${ }^{\#} \tilde{\mathcal{J}}_{b}(\Sigma)<+\infty$. Then there exist an integer $K \geqslant 0$ and an injection map $\phi: \mathbf{N}+K \mapsto$ $\mathcal{J}_{b}(\Sigma, 2) \times \mathbf{N}$ such that
(i) For any $k \in \mathbf{N}+K,[(\tau, x)] \in \mathcal{J}_{b}(\Sigma, 2)$ and $m \in \mathbf{N}$ satisfying $\phi(k)=([(\tau, x)]$, $m)$, there holds

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

where $x$ has minimal period $\tau$.
(ii) For any $k_{j} \in \mathbf{N}+K$, $k_{1}<k_{2}$, $\left(\tau_{j}, x_{j}\right) \in \mathcal{J}_{b}(\Sigma, 2)$ 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 3.4. (See Lemma 7.2 of [14].) 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 $i(\gamma) \geqslant n$. Then we have

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

Lemma 3.5. (See Lemma 7.3 of [14].) For any $(\tau, x) \in \mathcal{J}_{b}(\Sigma, 2)$ and $m \in \mathbf{N}$, we have

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

Proof of Theorem 1.1. By Theorem 1.1 of [14] we have ${ }^{\#} \tilde{\mathcal{J}}_{b}(\Sigma) \geqslant\left[\frac{n}{2}\right]+1$ for $n \in \mathbf{N}$. So we only need to prove Theorem 1.1 for the case $n \geqslant 3$ and $n$ is odd. The method of the proof is similar as that of [14].

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

$$
\begin{align*}
& H_{\Sigma}(x)=H_{\Sigma}(-x), \\
& H_{\Sigma}^{\prime}(x)=-H_{\Sigma}^{\prime}(-x), \\
& H_{\Sigma}^{\prime \prime}(x)=H_{\Sigma}^{\prime \prime}(-x) . \tag{3.74}
\end{align*}
$$

So $(\tau,-x) \in \mathcal{J}_{b}(\Sigma, 2)$. By (3.74) and the definition of $\gamma_{x}$ we have that

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

So we have

$$
\begin{align*}
& \left(i_{L_{0}}(x, m), v_{L_{0}}(x, m)\right)=\left(i_{L_{0}}(-x, m), v_{L_{0}}(-x, m)\right), \\
& \left(i_{L_{1}}(x, m), v_{L_{1}}(x, m)\right)=\left(i_{L_{1}}(-x, m), v_{L_{1}}(-x, m)\right), \quad \forall m \in \mathbf{N} . \tag{3.75}
\end{align*}
$$

So we can write

$$
\begin{equation*}
\tilde{\mathcal{J}}_{b}(\Sigma, 2)=\left\{\left[\left(\tau_{j}, x_{j}\right)\right] \mid j=1, \ldots, p\right\} \cup\left\{\left[\left(\tau_{k}, x_{k}\right)\right],\left[\left(\tau_{k},-x_{k}\right)\right] \mid k=p+1, \ldots, p+q\right\}, \tag{3.76}
\end{equation*}
$$

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

By Lemma 3.3 we have an integer $K \geqslant 0$ and an injection map $\phi: \mathbf{N}+K \rightarrow \mathcal{J}_{b}(\Sigma, 2) \times \mathbf{N}$. By (3.75), $\left(\tau_{k}, x_{k}\right)$ and ( $\tau_{k},-x_{k}$ ) have the same ( $i_{L_{0}}, v_{L_{0}}$ )-indices. So by Lemma 3.3, without loss of generality, we can further require that

$$
\begin{equation*}
\operatorname{Im}(\phi) \subseteq\left\{\left[\left(\tau_{k}, x_{k}\right)\right] \mid k=1,2, \ldots, p+q\right\} \times \mathbf{N} . \tag{3.77}
\end{equation*}
$$

By the strict convexity of $H_{\Sigma}$ and (6.19) of [14], we have

$$
\hat{i}_{L_{0}}\left(x_{k}\right)>0, \quad k=1,2, \ldots, p+q
$$

Applying Lemma 3.1 to the following associated symplectic paths

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

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

$$
\begin{align*}
& i_{L_{0}}\left(x_{k}, 2 m_{k}+1\right)=R+i_{L_{0}}\left(x_{k}\right)  \tag{3.78}\\
& i_{L_{0}}\left(x_{k}, 2 m_{k}-1\right)+v_{L_{0}}\left(x_{k}, 2 m_{k}-1\right)=R-\left(i_{L_{1}}\left(x_{k}\right)+n+S_{M_{k}}^{+}(1)-v_{L_{0}}\left(x_{k}\right)\right) \tag{3.79}
\end{align*}
$$

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

$$
\begin{align*}
& i_{L_{0}}\left(x_{k}, 4 m_{k}+2\right)=R+i_{L_{0}}\left(x_{k}, 2\right),  \tag{3.80}\\
& i_{L_{0}}\left(x_{k}, 4 m_{k}-2\right)+v_{L_{0}}\left(x_{k}, 4 m_{k}-2\right)=R-\left(i_{L_{1}}\left(x_{k}, 2\right)+n+S_{M_{k}}^{+}(1)-v_{L_{0}}\left(x_{k}, 2\right)\right) \tag{3.81}
\end{align*}
$$

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

$$
\begin{align*}
& i\left(x_{k}, 2 m_{k}+1\right)=2 R+i\left(x_{k}\right),  \tag{3.82}\\
& i\left(x_{k}, 2 m_{k}-1\right)+v\left(x_{k}, 2 m_{k}-1\right)=2 R-\left(i\left(x_{k}\right)+2 S_{M_{k}}^{+}(1)-v\left(x_{k}\right)\right), \tag{3.83}
\end{align*}
$$

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

$$
\begin{align*}
& i\left(x_{k}, 4 m_{k}+2\right)=2 R+i\left(x_{k}, 2\right),  \tag{3.84}\\
& i\left(x_{k}, 4 m_{k}-2\right)+v\left(x_{k}, 4 m_{k}-2\right)=2 R-\left(i\left(x_{k}, 2\right)+2 S_{M_{k}}^{+}(1)-v\left(x_{k}, 2\right)\right), \tag{3.85}
\end{align*}
$$

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 (3.77), we can set

$$
\phi(R-(s-1))=\left(\left[\left(\tau_{k(s)}, x_{k(s)}\right)\right], m(s)\right), \quad \forall s \in S:=\left\{1,2, \ldots,\left[\frac{n+1}{2}\right]+1\right\},
$$

where $k(s) \in\{1,2, \ldots, p+q\}$ and $m(s) \in \mathbf{N}$.
We continue our proof to study the symmetric and asymmetric orbits separately. Let

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

We shall prove that ${ }^{\#} S_{1} \leqslant p$ and ${ }^{\#} S_{2} \leqslant 2 q$, together with the definitions of $S_{1}$ and $S_{2}$, these yield Theorem 1.1.
Claim 1. ${ }^{\#} S_{1} \leqslant p$.
Proof. By the definition of $S_{1},\left(\left[\left(\tau_{k(s)}, x_{k(s)}\right)\right], m(s)\right)$ is symmetric when $k(s) \leqslant p$. We further prove that $m(s)=$ $2 m_{k(s)}$ for $s \in S_{1}$.

In fact, by the definition of $\phi$ and Lemma 3.3, for all $s=1,2, \ldots,\left[\frac{n+1}{2}\right]+1$ we have

$$
\begin{align*}
i_{L_{0}}\left(x_{k(s)}, m(s)\right) & \leqslant(R-(s-1))-1=R-s \\
& \leqslant i_{L_{0}}\left(x_{k(s)}, m(s)\right)+v_{L_{0}}\left(x_{k(s)}, m(s)\right)-1 . \tag{3.86}
\end{align*}
$$

By the strict convexity of $H_{\Sigma}$ and Lemma 2.2, we have $i_{L_{0}}\left(x_{k(s)}\right) \geqslant 0$, so there holds

$$
\begin{equation*}
i_{L_{0}}\left(x_{k(s)}, m(s)\right) \leqslant R-s<R \leqslant R+i_{L_{0}}\left(x_{k(s)}\right)=i_{L_{0}}\left(x_{k(s)}, 2 m_{k(s)}+1\right), \tag{3.87}
\end{equation*}
$$

for every $s=1,2, \ldots,\left[\frac{n+1}{2}\right]+1$, where we have used (3.78) in the last equality. Note that the proofs of (3.86) and (3.87) do not depend on the condition $s \in S_{1}$.

By Lemma 3.2, $\gamma_{x_{k}}$ satisfies conditions of Theorem 3.1 with $\tau=\frac{\tau_{k}}{2}$. 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 3.1 we have

$$
\begin{equation*}
i_{L_{1}}\left(x_{k}\right)+S_{M_{k}}^{+}(1)-v_{L_{0}}\left(x_{k}\right)>\frac{1-n}{2}, \quad \forall k=1, \ldots, p \tag{3.88}
\end{equation*}
$$

Also for $1 \leqslant s \leqslant\left[\frac{n+1}{2}\right]+1$, we have

$$
\begin{equation*}
-\frac{n+3}{2}=-\left(\left[\frac{n+1}{2}\right]+1\right) \leqslant-s . \tag{3.89}
\end{equation*}
$$

Hence by (3.86), (3.88) and (3.89), if $k(s) \leqslant p$ we have

$$
\begin{align*}
& i_{L_{0}}\left(x_{k(s)}, 2 m_{k(s)}-1\right)+v_{L_{0}}\left(x_{k(s)}, 2 m_{k(s)}-1\right)-1 \\
& \quad=R-\left(i_{L_{1}}\left(x_{k(s)}\right)+n+S_{M_{k(s)}}^{+}(1)-v_{L_{0}}\left(x_{k(s)}\right)\right)-1 \\
& \quad<R-\frac{1-n}{2}-1-n=R-\frac{n+3}{2} \leqslant R-s \\
& \quad \leqslant i_{L_{0}}\left(x_{k(s)}, m(s)\right)+v_{L_{0}}\left(x_{k(s)}, m(s)\right)-1 . \tag{3.90}
\end{align*}
$$

Thus by (3.87) and (3.90) and Lemma 3.5 of [14] we have

$$
\begin{equation*}
2 m_{k(s)}-1<m(s)<2 m_{k(s)}+1 . \tag{3.91}
\end{equation*}
$$

Hence

$$
\begin{equation*}
m(s)=2 m_{k(s)} . \tag{3.92}
\end{equation*}
$$

So we have

$$
\begin{equation*}
\phi(R-s+1)=\left(\left[\left(\tau_{k(s)}, x_{k(s)}\right)\right], 2 m_{k(s)}\right), \quad \forall s \in S_{1} . \tag{3.93}
\end{equation*}
$$

Then by the injectivity of $\phi$, it induces another injection map

$$
\begin{equation*}
\phi_{1}: S_{1} \rightarrow\{1, \ldots, p\}, \quad s \mapsto k(s) \tag{3.94}
\end{equation*}
$$

Therefore ${ }^{\#} S_{1} \leqslant p$. Claim 1 is proved.
Claim 2. ${ }^{\#} S_{2} \leqslant 2 q$.
Proof. By the formulas (3.82)-(3.85), and (59) of [13] (also Claim 4 on p. 352 of [16]), we have

$$
\begin{equation*}
m_{k}=2 m_{k+q} \quad \text { for } k=p+1, p+2, \ldots, p+q . \tag{3.95}
\end{equation*}
$$

We set $\mathcal{A}_{k}=i_{L_{1}}\left(x_{k}, 2\right)+S_{M_{k}}^{+}(1)-v_{L_{0}}\left(x_{k}, 2\right)$ and $\mathcal{B}_{k}=i_{L_{0}}\left(x_{k}, 2\right)+S_{M_{k}}^{+}(1)-v_{L_{1}}\left(x_{k}, 2\right), p+1 \leqslant k \leqslant p+q$, where $M_{k}=\gamma_{k}\left(2 \tau_{k}\right)=\gamma\left(\tau_{k}\right)^{2}$. By (3.7), we have

$$
\begin{equation*}
\mathcal{A}_{k}+\mathcal{B}_{k}=i\left(x_{k}, 2\right)+2 S_{M_{k}}^{+}(1)-v\left(x_{k}, 2\right)-n, \quad p+1 \leqslant k \leqslant p+q . \tag{3.96}
\end{equation*}
$$

By similar discussion of the proof of Lemma 3.2, for any $p+1 \leqslant k \leqslant p+q$ there exist $P_{k} \in \operatorname{Sp}(2 n)$ and $\tilde{M}_{k} \in$ $\operatorname{Sp}(2 n-2)$ such that

$$
\gamma\left(\tau_{k}\right)=P_{k}^{-1}\left(I_{2} \diamond \tilde{M}_{k}\right) P_{k} .
$$

Hence by Lemma 3.4 and (3.96), we have

$$
\begin{equation*}
\mathcal{A}_{k}+\mathcal{B}_{k} \geqslant n+2-n=2 . \tag{3.97}
\end{equation*}
$$

By Theorem 2.1, there holds

$$
\begin{equation*}
\left|\mathcal{A}_{k}-\mathcal{B}_{k}\right|=\left|\left(i_{L_{0}}\left(x_{k}, 2\right)+v_{L_{0}}\left(x_{k}, 2\right)\right)-\left(i_{L_{1}}\left(x_{k}, 2\right)+v_{L_{1}}\left(x_{k}, 2\right)\right)\right| \leqslant n . \tag{3.98}
\end{equation*}
$$

So by (3.97) and (3.98) we have

$$
\begin{equation*}
\mathcal{A}_{k} \geqslant \frac{1}{2}\left(\left(\mathcal{A}_{k}+\mathcal{B}_{k}\right)-\left|\mathcal{A}_{k}-\mathcal{B}_{k}\right|\right) \geqslant \frac{2-n}{2}, \quad p+1 \leqslant k \leqslant p+q . \tag{3.99}
\end{equation*}
$$

By (3.81), (3.86), (3.89), (3.95) and (3.99), for $p+1 \leqslant k(s) \leqslant p+q$ we have

$$
\begin{align*}
& i_{L_{0}}\left(x_{k(s)}, 2 m_{k(s)}-2\right)+v_{L_{0}}\left(x_{k(s)}, 2 m_{k(s)}-2\right)-1 \\
& \quad=i_{L_{0}}\left(x_{k(s)}, 4 m_{k(s)+q}-2\right)+v_{L_{0}}\left(x_{k(s)}, 4 m_{k(s)+q}-2\right)-1 \\
& \quad=R-\left(i_{L_{1}}\left(x_{k(s)}, 2\right)+n+S_{M_{k(s)}}^{+}(1)-v_{L_{0}}\left(x_{k(s)}, 2\right)\right)-1 \\
& \quad=R-\mathcal{A}_{k(s)}-1-n \\
& \quad \leqslant R-\frac{2-n}{2}-1-n \\
& \quad=R-\left(2+\frac{n}{2}\right) \\
& \quad<R-\frac{n+3}{2} \\
& \quad \leqslant R-s \\
& \quad \leqslant i_{L_{0}}\left(x_{k(s)}, m(s)\right)+v_{L_{0}}\left(x_{k(s)}, m(s)\right)-1 . \tag{3.100}
\end{align*}
$$

Thus by (3.87), (3.100) and Lemma 3.5, we have

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

So

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

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

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

Thus by the injectivity of the map $\phi$ from Lemma 3.3, we have

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

which yields Claim 2.
By Claim 1 and Claim 2, we have

$$
{ }^{\#} \tilde{\mathcal{J}}_{b}(\Sigma)={ }^{\#} \tilde{\mathcal{J}}_{b}(\Sigma, 2)=p+2 q \geqslant{ }^{\#} S_{1}+{ }^{\#} S_{2}=\left[\frac{n+1}{2}\right]+1 .
$$

The proof of Theorem 1.1 is complete.
Proof of Theorem 1.2. By [13], there are at least $n$ closed characteristics on every $C^{2}$ compact convex central symmetric hypersurface $\Sigma$ of $\mathbf{R}^{2 n}$. Hence by Example 1.1 the assumption of Theorem 1.2 is reasonable. Here we prove the case $n=5$, the proof of the case $n=4$ is the same.

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 [22], a closed characteristic $x$ on $\Sigma$ can became a dual brake orbit after suitable time translation if and only if $x(\mathbf{R})=-N x(\mathbf{R})$. So by Lemma 3.1 of [22] again, if a closed characteristic $x$ on $\Sigma$ can both became brake orbits and dual brake orbits after suitable translation, then $x(\mathbf{R})=N x(\mathbf{R})=-N x(\mathbf{R})$. Thus $x(\mathbf{R})=-x(\mathbf{R})$.

Since we also have $-N \Sigma=\Sigma,(-N)^{2}=I_{2 n}$ and $(-N) J=-J(-N)$, dually by the same proof of Theorem 1.1, there are at least $[(n+1) / 2]+1=4$ geometrically distinct dual brake orbits on $\Sigma$.

If there are exactly 5 closed characteristics on $\Sigma$. By Theorem 1.1, four closed characteristics of them must be brake orbits after suitable time translation, then the fifth, say $y$, must be brake orbits after suitable time translation, otherwise $N y(-\cdot)$ will be the sixth geometrically distinct closed characteristic on $\Sigma$ which yields a contradiction. Hence all closed characteristics on $\Sigma$ must be brake orbits on $\Sigma$. By the same argument we can prove that all closed characteristics on $\Sigma$ must be dual brake orbits on $\Sigma$. Then by the argument in the second paragraph of the proof of this theorem, all these five closed characteristics on $\Sigma$ must be symmetric. Hence all of them must be symmetric brake orbits after suitable time translation. Thus we have proved the case $n=5$ of Theorem 1.2 and the proof of Theorem 1.2 is complete.

## References

[1] A. Ambrosetti, V. Benci, Y. Long, A note on the existence of multiple brake orbits, Nonlinear Anal. 21 (1993) 643-649.
[2] V. Benci, Closed geodesics for the Jacobi metric and periodic solutions of prescribed energy of natural Hamiltonian systems, Ann. Inst. H. Poincaré Anal. Non Linéaire 1 (1984) 401-412.
[3] V. Benci, F. Giannoni, A new proof of the existence of a brake orbit, in: Advanced Topics in the Theory of Dynamical Systems, in: Notes Rep. Math. Sci. Eng., vol. 6, 1989, pp. 37-49.
[4] S. Bolotin, Libration motions of natural dynamical systems, Vestnik Moskov. Univ. Ser. I Mat. Mekh. 6 (1978) $72-77$ (in Russian).
[5] S. Bolotin, V.V. Kozlov, Librations with many degrees of freedom, J. Appl. Math. Mech. 42 (1978) 245-250 (in Russian).
[6] S.E. Cappell, R. Lee, E.Y. Miller, On the Maslov-type index, Comm. Pure Appl. Math. 47 (1994) 121-186.
[7] I. Ekeland, Convexity Methods in Hamiltonian Mechanics, Springer-Verlag, Berlin, 1990.
[8] H. Gluck, W. Ziller, Existence of periodic solutions of conservative systems, in: Seminar on Minimal Submanifolds, Princeton University Press, 1983, pp. 65-98.
[9] E.W.C. van Groesen, Analytical mini-max methods for Hamiltonian brake orbits of prescribed energy, J. Math. Anal. Appl. 132 (1988) 1-12.
[10] K. Hayashi, Periodic solution of classical Hamiltonian systems, Tokyo J. Math. 6 (1983) 473-486.
[11] C. Liu, Maslov-type index theory for symplectic paths with Lagrangian boundary conditions, Adv. Nonlinear Stud. 7 (1) (2007) 131-161.
[12] C. Liu, Asymptotically linear Hamiltonian systems with Lagrangian boundary conditions, Pacific J. Math. 232 (1) (2007) 233-255.
[13] C. Liu, Y. Long, C. Zhu, Multiplicity of closed characteristics on symmetric convex hypersurfaces in $\mathbf{R}^{2 n}$, Math. Ann. 323 (2) (2002) 201-215.
[14] C. Liu, D. Zhang, Iteration theory of $L$-index and multiplicity of brake orbits, arXiv:0908.0021v1 [math.SG].
[15] Y. Long, Bott formula of the Maslov-type index theory, Pacific J. Math. 187 (1999) 113-149.
[16] Y. Long, Index Theory for Symplectic Paths with Applications, Birkhäuser, Basel, 2002.
[17] Y. Long, D. Zhang, C. Zhu, Multiple brake orbits in bounded convex symmetric domains, Adv. Math. 203 (2006) 568-635.
[18] Y. Long, C. Zhu, Closed characteristics on compact convex hypersurfaces in $\mathbf{R}^{2 n}$, Ann. of Math. 155 (2002) 317-368.
[19] P.H. Rabinowitz, On the existence of periodic solutions for a class of symmetric Hamiltonian systems, Nonlinear Anal. 11 (1987) 599-611.
[20] H. Seifert, Periodische Bewegungen mechanischer Systeme, Math. Z. 51 (1948) 197-216.
[21] A. Szulkin, An index theory and existence of multiple brake orbits for star-shaped Hamiltonian systems, Math. Ann. 283 (1989) 241-255.
[22] D. Zhang, Brake type closed characteristics on reversible compact convex hypersurfaces in $\mathbf{R}^{2 n}$, Nonlinear Anal. 74 (2011) 3149-3158.
[23] D. Zhang, Minimal period problems for brake orbits of nonlinear autonomous reversible semipositive Hamiltonian systems, Discrete Contin. Dyn. Syst. (2013), in press, arXiv:1110.6915v1 [math.SG].


[^0]:    * Corresponding author.

    E-mail addresses: zhangdz@nankai.edu.cn (D. Zhang), liucg@nankai.edu.cn (C. Liu).
    1 Partially supported by the NSF of China (10801078, 11171314, 11271200) and Nankai University.
    2 Partially supported by the NSF of China (11071127, 10621101), 973 Program of MOST (2011CB808002) and SRFDP.

