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# Singularly perturbed spectral problems with Neumann boundary conditions 

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The paper deals with the Neumann spectral problem for a singularly perturbed second-order elliptic operator with bounded lower order terms. The main goal is to provide a refined description of the limit behaviour of the principal eigenvalue and eigenfunction. Using the logarithmic transformation, we reduce the studied problem to an additive eigenvalue problem for a singularly perturbed Hamilton-Jacobi equation. Then assuming that the Aubry set of the Hamiltonian consists of a finite number of points or limit cycles situated in the domain or on its boundary, we find the limit of the eigenvalue and formulate the selection criterion that allows us to choose a solution of the limit Hamilton-Jacobi equation which gives the logarithmic asymptotics of the principal eigenfunction.

Keywords: singularly perturbed operator; Neumann spectral problem; HamiltonJacobi equation

AMS Subject Classifications: 35B27; 35B40

## 1. Introduction

This paper is devoted to the asymptotic analysis of the first eigenpair for singularly perturbed spectral problem, depending on the small parameter $\varepsilon>0$, for the elliptic equation

$$
\begin{equation*}
\varepsilon a_{i j}(x) \frac{\partial^{2} u_{\varepsilon}}{\partial x_{i} \partial x_{j}}+b_{i}(x) \frac{\partial u_{\varepsilon}}{\partial x_{i}}+c(x) u_{\varepsilon}=\lambda_{\varepsilon} u \tag{1.1}
\end{equation*}
$$

in a smooth bounded domain $\Omega \subset \mathbb{R}^{N}$ with the boundary condition

$$
\begin{equation*}
\frac{\partial u_{\varepsilon}}{\partial v}=0 \tag{1.2}
\end{equation*}
$$

on $\partial \Omega$, where $\frac{\partial}{\partial \nu}$ denotes derivative with respect to the external normal.

[^0]The bottom of the spectrum of elliptic operators plays a crucial role in many applications. In particular, the first eigenvalue and the corresponding eigenfunction of (1.1)-(1.2), are important in understanding the large-time behaviour of the underlying non-stationary convection-diffusion model with reflecting boundary. Due to the Krein-Rutman theorem, the first eigenvalue $\lambda_{\varepsilon}$ of (1.1)-(1.2) (the eigenvalue with the maximal real part) is simple and real, the corresponding eigenfunction $u_{\varepsilon}$ can be chosen to satisfy $u_{\varepsilon}(x)>0$ in $\Omega$.

The goal of this work is to study the asymptotic behaviour of $\lambda_{\varepsilon}$ and $u_{\varepsilon}$ as $\varepsilon \rightarrow 0$. While in the case of constant function $c(x)$ in (1.1), the first eigenpair is (trivially) explicitly found, the asymptotic behaviour of the first eigenpair is quite non-trivial when $c(x)$ is a non-constant function, in particular, the eigenfunction might exhibit an exponential localization.

Boundary value problems for singularly perturbed elliptic operators have been actively studied starting from 1950s. We mention here a pioneering work,[1] where for a wide class of operators (so-called regularly degenerated operators), the asymptotics of solutions were obtained.

In the works [2-4] (see also [5]), the principal eigenvalue of singularly perturbed convection-diffusion equations with the Dirichlet boundary condition was investigated by means of large deviation techniques for diffusion processes with small diffusion. In [6], the estimates for the principal eigenvalue were obtained by comparison arguments and elliptic techniques.

The case when convection vector field has a finite number of hyperbolic equilibrium points and cycles was studied in [7] where methods of dynamical systems are combined with those of stochastic differential equations. These results were generalized in [8] to the case when the boundary of domain is invariant with respect to convection vector field. Similar problem in the presence of zero-order term was considered in [9].

In [10], the viscosity solution techniques for singularly perturbed Hamilton-Jacobi equation were used in order to study the principal eigenfunction of the adjoint Neumann convection-diffusion problem. The logarithmic asymptotics of the eigenfunction were constructed.

The work [11] deals with the principal eigenpair of operators with a large zero-order term on a compact Riemannian manifold. The approach developed in this work is based on large deviation and variational techniques.

Dirichlet spectral problem for singularly perturbed operators with rapidly oscillating locally periodic coefficients was studied in [12] and [13]. In [12] with the help of viscosity solution method, the limit of the principal eigenvalue and the logarithmic asymptotics of the principal eigenfunction were found. These asymptotics were improved in [12] and [13] using the blow up analysis.

In the present work, when studying problem (1.1)-(1.2), we make use of the standard viscosity solution techniques in order to obtain the logarithmic asymptotics of the principal eigenfunction. However, the limit Hamilton-Jacobi equation in general is not uniquely solvable and does not give information about the limit behaviour of $\lambda_{\varepsilon}$. Therefore, we have to consider higher order approximations in (1.1)-(1.2). Under rather general assumptions on the structure of the Aubry set of the limit Hamiltonian, we find the limit of $\lambda_{\varepsilon}$ and can choose the solution of the limit problem which determines the asymptotics of the principal eigenfunction. Notice that we did not succeed to make the blow up analysis work in the case under consideration. In this case, for components of the Aubry set located on the boundary, the natural rescaling still leads to singularly perturbed operators. Instead, we study a refined structure of solutions of the limit Hamilton-Jacobi equation in the vicinity of the Aubry set.

This allows us to construct test functions that satisfy the perturbed equation up to higher order.

We also would like to remark that, with obvious modifications, the results of this work as well as the developed techniques remain valid for the boundary condition of the form

$$
\frac{\partial u_{\varepsilon}}{\partial \beta}=0
$$

where $\beta$ is a $C^{2}$-smooth vector field on $\partial \Omega$ non-tangential at any point of $\partial \Omega$. In particular, conormal vector field $\beta_{i}=a_{i j} v_{j}$ can be considered.

## 2. Problem set-up and results

We study problem (1.1)-(1.2) under the following assumptions on the operator coefficients and the domain:
(a1) $\Omega$ is a bounded domain in $\mathbb{R}^{N}, N \geq 2$, with a $C^{2}$ boundary;
(a2) all the coefficients are $C^{2}$-functions in $\bar{\Omega}$;
(a3) the matrix $\left(a_{i j}\right)$ is symmetric and uniformly elliptic.
Further assumptions on the vector field $b$ will be formulated later on.
Since $u_{\varepsilon}>0$ in $\Omega$, we can represent $u_{\varepsilon}$ in the form

$$
u_{\varepsilon}=e^{-W_{\varepsilon}(x) / \varepsilon},
$$

this results in the following non-linear PDE

$$
\begin{equation*}
-a_{i j}(x) \frac{\partial^{2} W_{\varepsilon}}{\partial x_{i} \partial x_{j}}+\frac{1}{\varepsilon} H\left(\nabla W_{\varepsilon}, x\right)+c(x)=\lambda_{\varepsilon} \quad \text { in } \Omega \tag{2.1}
\end{equation*}
$$

or

$$
\begin{equation*}
-\varepsilon a_{i j}(x) \frac{\partial^{2} W_{\varepsilon}}{\partial x_{i} \partial x_{j}}+H\left(\nabla W_{\varepsilon}, x\right)+\varepsilon c(x)=\varepsilon \lambda_{\varepsilon} \quad \text { in } \Omega \tag{2.2}
\end{equation*}
$$

with the boundary condition

$$
\begin{equation*}
\frac{\partial W_{\varepsilon}}{\partial \nu}=0 \quad \text { on } \partial \Omega \tag{2.3}
\end{equation*}
$$

where

$$
\begin{equation*}
H(p, x)=a_{i j}(x) p_{i} p_{j}-b_{i}(x) p_{i} \tag{2.4}
\end{equation*}
$$

is a function to be referred to as a Hamiltonian. Passing to the limit as $\varepsilon \rightarrow 0$ in (2.2), with the help of the standard approach based on the maximum principle, we can show that $W_{\varepsilon}$ converges uniformly (up to extracting a subsequence) to a viscosity solution $W(x)$ of the Hamiltion-Jacobi equation

$$
\begin{equation*}
H(\nabla W(x), x)=0 \quad \text { in } \Omega \tag{2.5}
\end{equation*}
$$

with the boundary condition

$$
\begin{equation*}
\frac{\partial W}{\partial \nu}=0 \quad \text { on } \partial \Omega \tag{2.6}
\end{equation*}
$$

Recall that a function $W \in C(\bar{\Omega})$ is called a viscosity solution of Equation (2.5) if for every test function $\Phi \in C^{\infty}(\bar{\Omega})$, the following holds:

- if $W-\Phi$ attains a maximum at a point $\xi \in \Omega$, then $W(\nabla \Phi(\xi), \xi) \leq 0$;
- if $W-\Phi$ attains a minimum at $\xi \in \Omega$, then $W(\nabla \Phi(\xi), \xi) \geq 0$.

The boundary condition (2.6) is understood in the following sense, $\forall \Phi \in C^{\infty}(\bar{\Omega})$

- if $W-\Phi$ attains a maximum at $\xi \in \partial \Omega$, then $\min \left\{H(\nabla \Phi(\xi), \xi), \frac{\partial \Phi}{\partial v}(\xi)\right\} \leq 0$;
- if $W-\Phi$ attains a minimum at $\xi \in \partial \Omega$, then $\max \left\{H(\nabla \Phi(\xi), \xi), \frac{\partial \Phi}{\partial \nu}(\xi)\right\} \geq 0$.

It is known [14] that every solution of problem (2.5)-(2.6) has the representation

$$
\begin{equation*}
W(x)=\inf _{y \in \mathcal{A}_{H}}\left\{d_{H}(x, y)+W(y)\right\}, \tag{2.7}
\end{equation*}
$$

where $\mathcal{A}_{H}$ is so-called Aubry set and $d_{H}(x, y)$ is a distance function. To define $\mathcal{A}_{H}$ and $d_{H}(x, y)$ consider solutions of the following Skorohod problem:

$$
\left\{\begin{array}{l}
\eta(t) \in \bar{\Omega}, \quad t \geq 0  \tag{2.8}\\
\dot{\eta}(t)+\alpha(t) v(\eta(t))=v(t) \quad \text { with } \alpha(t) \geq 0 \text { and } \alpha(t)=0 \text { when } \eta(t) \notin \partial \Omega \\
\eta(0)=x,
\end{array}\right.
$$

where $v \in L^{1}\left((0, \infty) ; \mathbb{R}^{N}\right)$ is a given vector field and $x \in \bar{\Omega}$ is a given initial point, while the curve $\eta \in W_{\text {loc }}^{1,1}\left((0, \infty) ; \mathbb{R}^{N}\right)$ and the function $\alpha \in L^{1}\left((0, \infty) ; \mathbb{R}_{+}\right)$are unknowns. Under our standing assumptions on $\Omega\left(\partial \Omega \in C^{2}\right)$, the Skorohod problem (2.8) has a solution, see [14].

Consider now the Legendre transform $L(v, x)=\sup _{p \in \mathbb{R}^{N}}(v \cdot p-H(p, x))$ and define the distance function

$$
\begin{equation*}
d_{H}(x, y)=\inf \left\{\int_{0}^{t} L(-v(s), \eta(s)) \mathrm{d} s, \quad \eta \text { solves }(2.8), \eta(0)=x, \eta(t)=y, t>0\right\} . \tag{2.9}
\end{equation*}
$$

Next, we recall the variational definition of the Aubry set

$$
\begin{equation*}
x \in \mathcal{A}_{H} \Longleftrightarrow \forall \delta>0 \quad \inf \left\{\int_{0}^{t} L(-v(s), \eta(s)) \mathrm{d} s, \quad \eta \text { solves }(2.8), \eta(0)=\eta(t)=x, t>\delta\right\}=0 . \tag{2.10}
\end{equation*}
$$

In this work, we assume that the Aubry set has a finite number of connected components

$$
\begin{align*}
& \mathcal{A}_{H}=\bigcup_{\text {finite }} \mathcal{A}_{k} \text { and each } \mathcal{A}_{k} \text { is either an isolated point } \\
& \quad \text { or a closed curve lying entirely either in } \Omega \text { or on } \partial \Omega . \tag{2.11}
\end{align*}
$$

Additionally, we assume that

$$
\text { if } \mathcal{A}_{k} \subset \Omega \text { then } \mathcal{A}_{k} \text { is either a hyperbolic fixed point }, ~ \begin{align*}
& \text { or a hyperbolic limit cycle of the } O D E \dot{x}=b(x) ;
\end{align*}
$$

if $\mathcal{A}_{k} \subset \partial \Omega$ then the normal component $b_{\nu}(x)$ of the field $b(x)$ is strictly positive on $\mathcal{A}_{k}$ and $\mathcal{A}_{k}$ is either a hyperbolic fixed point or a hyperbolic limit cycle of the ODE $\dot{x}=b_{\tau}(x)$ on $\partial \Omega$, where $b_{\tau}(x)$ denotes the tangential component of $b(x)$ on $\partial \Omega$.

Remark 1 Note that the Aubry set $\mathcal{A}_{H}$ does not depend on the coefficients $a_{i j}(x)$, it is completely defined by the drift field $b(x)$. This fact follows from the variational definition (2.10) of the Aubry set observing that the Lagrangian $L(v, x)$ is given by $L(v, x)=\frac{1}{4} a^{i j}(x)\left(v_{i}+\right.$ $\left.b_{i}(x)\right)\left(v_{j}+b_{j}(x)\right)$, where $\left(a^{i j}(x)\right)_{i, j=\overline{1, N}}$ is the matrix inverse to $\left(a_{i j}(x)\right)_{i, j=\overline{1, N}}$. More specifically, $\mathcal{A}_{H}$ is determined by the dynamical system $\mathcal{S}$ corresponding to the Skorohod problem (2.8) with $v(t)=b(\eta(t))$ and conditions (2.11)-(2.13) require that the $\omega$-limit set of $\mathcal{S}$ consists of a finite number of hyperbolic fixed points or limit cycles. One observes that the $\omega$-limit set of $\mathcal{S}$ is always non-empty. Furthermore, in the case of general position, the Aubry set $\mathcal{A}_{H}$ consists of a finite number of hyperbolic fixed points and limit cycles of $\mathcal{S}$.

In order to state the main result of this work, we assign to each component $\mathcal{A}_{k}$ of $\mathcal{A}_{H}$ a number $\sigma\left(\mathcal{A}_{k}\right)$ as follows. If $\mathcal{A}_{k}$ is a fixed point $\{\xi\}$ of the $\mathrm{ODE} \dot{x}=b(x)$ and $\xi \in \Omega$, linearizing the ODE near $\xi$ to get $\dot{z}=B(\xi) z$ we define $\sigma\left(\mathcal{A}_{k}\right)$ by

$$
\begin{equation*}
\sigma\left(\mathcal{A}_{k}\right)=-\sum_{\theta_{i}>0} \theta_{i}+c(\xi) \tag{2.14}
\end{equation*}
$$

where $\theta_{i}$ are the real parts of eigenvalues of the matrix $B(\xi)$. Note that the hyperbolicity of the fixed point means that the eigenvalues of $B(\xi)$ cannot have zero real part. If $\mathcal{A}_{k}=\{\xi\}$ and $\xi \in \partial \Omega$, consider the $\mathrm{ODE} \dot{x}=b_{\tau}(x)$ on $\partial \Omega$ in a neighbourhood of the point $\xi$. Passing to the linearized $\mathrm{ODE} \dot{z}=B_{\tau}(\xi) z$ in the tangent plane to $\partial \Omega$ at the point $\xi$, we denote by $\tilde{\theta}_{i}$ the real parts of the eigenvalues of $B_{\tau}(\xi)$ and set

$$
\begin{equation*}
\sigma\left(\mathcal{A}_{k}\right)=-\sum_{\tilde{\theta}_{i}>0} \tilde{\theta}_{i}+c(\xi) \tag{2.15}
\end{equation*}
$$

Consider now the case when $\left\{\mathcal{A}_{k}\right\} \subset \Omega$ is a limit cycle of ODE $\dot{x}=b(x)$. Let $P>0$ be the minimal period of the cycle and let $\Theta_{i}$ be the absolute values of eigenvalues of the linearized Poincaré map. (Recall that the limit cycle is said hyperbolic if there are no eigenvalues of linearized Poincaré map with absolute value equal to 1.) We define now $\sigma\left(\mathcal{A}_{k}\right)$ by setting

$$
\begin{equation*}
\sigma\left(\mathcal{A}_{k}\right)=-\frac{1}{P} \sum_{\Theta_{i}>1} \log \Theta_{i}+\frac{1}{P} \int_{0}^{P} c(\xi(t)) \mathrm{d} t \tag{2.16}
\end{equation*}
$$

where $\xi(t)$ solves $\dot{\xi}=b(\xi)$ and $\xi(t) \in \mathcal{A}_{k}$.
Finally, in the case when $b_{v}>0$ on $\mathcal{A}_{k}$ and $\mathcal{A}_{k}$ is a limit cycle of the $\mathrm{ODE} \dot{x}=b_{\tau}(x)$ on $\partial \Omega$, we set

$$
\begin{equation*}
\sigma\left(\mathcal{A}_{k}\right)=-\frac{1}{P} \sum_{\widetilde{\Theta}_{i}>1} \log \widetilde{\Theta}_{i}+\frac{1}{P} \int_{0}^{P} c(\xi(t)) \mathrm{d} t \tag{2.17}
\end{equation*}
$$

where $\dot{\xi}=b_{\tau}(\xi)$ and $\xi(t) \in \mathcal{A}_{k}, P$ is the minimal period and $\widetilde{\Theta}_{i}$ are the absolute values of the eigenvalues of the linearized Poincaré map.

The main result of this work is
Theorem 2 Let conditions (a1)-(a3) be fulfilled, and assume that the Aubry set $\mathcal{A}_{H}$ satisfies (2.11), (2.12) and (2.13). Then the first eigenvalue $\lambda_{\varepsilon}$ of (1.1) converges as $\varepsilon \rightarrow 0$ to

$$
\begin{equation*}
\lim _{\varepsilon \rightarrow 0} \lambda_{\varepsilon}=\max \left\{\sigma\left(\mathcal{A}_{k}\right) ; \mathcal{A}_{k} \subset \mathcal{A}_{H}\right\} \tag{2.18}
\end{equation*}
$$

where $\sigma\left(\mathcal{A}_{k}\right)$ is given by (2.14) or (2.15) if $\mathcal{A}_{k}$ is a fixed point in $\Omega$ or on $\partial \Omega$, and $\sigma\left(\mathcal{A}_{k}\right)$ is defined by (2.16) or (2.17) if $\mathcal{A}_{k}$ is a limit cycle in $\Omega$ or on $\partial \Omega$. Moreover, if the maximum in (2.18) is attained at exactly one component $\mathcal{M}:=\mathcal{A}_{k_{0}}$, then the scaled logarithmic transform $w_{\varepsilon}=-\varepsilon \log u_{\varepsilon}$ of the first eigenfunction $u_{\varepsilon}$ (normalized by max $u_{\varepsilon}=1$ ) converges uniformly in $\bar{\Omega}$ to the maximal viscosity solution $W$ of (2.5) and (2.6) vanishing on $\mathcal{M}$, i.e. $W(x)=d_{H}(x, \mathcal{M})$.

## 3. Passing to the limit by vanishing viscosity techniques

In this section, we pass to the limit, as $\varepsilon \rightarrow 0$, in Equation (2.2) and boundary condition (2.3) to get problem (2.5) and (2.6). We use the standard technique based on the maximum principle and the a priori uniform $W^{1, \infty}$ bounds for $W_{\varepsilon}$ obtained by Bernstein's method (originally proposed in $[15,16]$ and further developed in, e.g. $[17,18]$ ).

First, considering (2.2) at the maximum and minimum points of $W_{\varepsilon}(x)$, we easily get
Lemma 3 The first eigenvalue $\lambda_{\varepsilon}$ satisfies the estimates $\min c(x) \leq \lambda_{\varepsilon} \leq \max c(x)$.
Next, we establish the $W^{1, \infty}$ bound for $W_{\varepsilon}$ in
Lemma 4 Let $u_{\varepsilon}$ be normalized by $\max u_{\varepsilon}=1$ (i.e. $\min W_{\varepsilon}=0$ ). Then $\left\|W_{\varepsilon}\right\|_{W^{1, \infty}(\Omega)} \leq$ $C$ with a constant $C$ independent of $\varepsilon$.

Proof Following [19] observe that the boundary condition $\frac{\partial W_{\varepsilon}}{\partial v}=0$ yields the pointwise bound

$$
\frac{\partial}{\partial \nu}\left|\nabla W_{\varepsilon}\right|^{2} \leq C\left|\nabla W_{\varepsilon}\right|^{2} \quad \text { on } \partial \Omega \text {. }
$$

Therefore, for an appropriate positive function $\phi \in C^{2}(\bar{\Omega})$,

$$
\begin{equation*}
\frac{\partial}{\partial \nu}\left(\phi\left|\nabla W_{\varepsilon}\right|^{2}\right) \leq-\left(\phi\left|\nabla W_{\varepsilon}\right|^{2}\right) \quad \text { on } \partial \Omega \tag{3.1}
\end{equation*}
$$

Next, we use Bernstein's method to obtain a uniform bound for $\omega_{\varepsilon}(x)=\phi(x)\left|\nabla W_{\varepsilon}(x)\right|^{2}$, following closely the line of [20], Lemma 1.2. In view of (3.1) either $\left|\nabla W_{\varepsilon}\right| \equiv 0$ and we have nothing to prove, or $\max w_{\varepsilon}$ is attained at a point $\xi \in \Omega$. In the latter case, we have $\nabla \omega_{\varepsilon}(\xi)=0$ and

$$
a_{i j} \frac{\partial^{2}}{\partial x_{i} \partial x_{j}}\left(\phi\left|\nabla W_{\varepsilon}\right|^{2}\right) \leq 0 \quad \text { at } x=\xi .
$$

Expanding the left-hand side of this inequality, we get

$$
\begin{align*}
2 \varepsilon \phi a_{i j} \frac{\partial^{2} W_{\varepsilon}}{\partial x_{i} \partial x_{k}} \frac{\partial^{2} W_{\varepsilon}}{\partial x_{j} \partial x_{k}} \leq & -2 \varepsilon \phi a_{i j} \frac{\partial^{3} W_{\varepsilon}}{\partial x_{i} \partial x_{j} \partial x_{k}} \frac{\partial W_{\varepsilon}}{\partial x_{k}} \\
& -4 \varepsilon a_{i j} \frac{\partial \phi}{\partial x_{j}} \frac{\partial^{2} W_{\varepsilon}}{\partial x_{i} \partial x_{k}} \frac{\partial W_{\varepsilon}}{\partial x_{k}}-\varepsilon a_{i j} \frac{\partial^{2} \phi}{\partial x_{i} \partial x_{j}}\left|\nabla W_{\varepsilon}\right|^{2} \tag{3.2}
\end{align*}
$$

Using (2.2) we obtain
$-\varepsilon \phi a_{i j} \frac{\partial^{3} W_{\varepsilon}}{\partial x_{i} \partial x_{j} \partial x_{k}} \frac{\partial W_{\varepsilon}}{\partial x_{k}} \leq \varepsilon \phi \frac{\partial a_{i j}}{\partial x_{k}} \frac{\partial^{2} W_{\varepsilon}}{\partial x_{i} \partial x_{j}} \frac{\partial W_{\varepsilon}}{\partial x_{k}}+C\left(\omega_{\varepsilon}^{3 / 2}+\omega_{\varepsilon}+\omega_{\varepsilon}^{1 / 2}+1\right), \quad$ at $x=\xi$,
where we have also exploited the fact that $\nabla \omega_{\varepsilon}(\xi)=0$. Substitute now (3.3) into (3.2) to derive

$$
\begin{equation*}
\varepsilon \frac{\partial^{2} W_{\varepsilon}}{\partial x_{i} \partial x_{k}} \frac{\partial^{2} W_{\varepsilon}}{\partial x_{i} \partial x_{k}} \leq C\left(\omega_{\varepsilon}^{3 / 2}+1\right), \quad \text { at } x=\xi \tag{3.4}
\end{equation*}
$$

On the other hand, it follows from (2.2) that

$$
\begin{equation*}
\omega_{\varepsilon} \leq C\left(\varepsilon \sum\left|\frac{\partial^{2} W_{\varepsilon}}{\partial x_{i} \partial x_{j}}\right|+1\right) \tag{3.5}
\end{equation*}
$$

Combining (3.4) and (3.5), we obtain $\omega_{\varepsilon} \leq C$ and the required uniform bound follows.
With a priori bounds from Lemma 3 and 4, it is quite standard to pass to the limit in (2.2). Indeed, up to extracting a subsequence, $W_{\varepsilon} \rightarrow W$ uniformly in $\bar{\Omega}$ and $\lambda_{\varepsilon} \rightarrow \lambda$. Consider a test function $\Phi \in C^{\infty}(\bar{\Omega})$ and assume that $W-\Phi$ attains strict maximum at a point $\xi$. Then, $W_{\varepsilon}-\Phi$ attains local maximum at $\xi_{\varepsilon}$ such that $\xi_{\varepsilon} \rightarrow \xi$ as $\varepsilon \rightarrow 0$. If $\xi_{\varepsilon} \in \Omega$, then $\nabla W_{\varepsilon}\left(\xi_{\varepsilon}\right)=\nabla \Phi\left(\xi_{\varepsilon}\right)$ and

$$
a_{i j} \frac{\partial^{2} W_{\varepsilon}}{\partial x_{i} \partial x_{j}}\left(\xi_{\varepsilon}\right) \leq a_{i j} \frac{\partial^{2} \Phi}{\partial x_{i} \partial x_{j}}\left(\xi_{\varepsilon}\right) \quad \text { if } \xi_{\varepsilon} \in \Omega
$$

and $\frac{\partial \Phi}{\partial \nu}\left(\xi_{\varepsilon}\right) \leq 0$ if $\xi_{\varepsilon} \in \partial \Omega$. Passing to the limit as $\varepsilon \rightarrow 0$ and using (2.2) and Lemma 3, we get

$$
H(\nabla \Phi(\xi), \xi) \leq 0 \text { if } \xi \in \Omega, \quad \text { and } \min \left\{H(\nabla \Phi(\xi), \xi), \frac{\partial \Phi}{\partial \nu}(\xi)\right\} \leq 0 \text { if } \xi \in \partial \Omega
$$

Arguing similarly in the case when $\xi$ is a strict minimum of $W-\Phi$, we conclude that $W$ is a viscosity solution of (2.5) and (2.6).

## 4. Matching lower and upper bounds for eigenvalues and selection of the solution of (2.5)-(2.6)

Due to the results of Section 3, we can assume, passing to a subsequence if necessary, that eigenvalues $\lambda_{\varepsilon}$ converge to a finite limit $\lambda$ and functions $W_{\varepsilon}$ converge uniformly in $\bar{\Omega}$ to a solution $W$ of problem (2.5) and (2.6) as $\varepsilon \rightarrow 0$. In the following four steps, we prove that $\lambda$ and $W(x)$ are described by Theorem 2.

Step I: Significant component(s) of $\mathcal{A}_{H}$. Recall the definition of the partial order relation $\preceq$ on $\mathcal{A}_{H}$ introduced in [12] as follows:

$$
\begin{equation*}
\mathcal{A}^{\prime} \preceq \mathcal{A}^{\prime \prime} \Longleftrightarrow W\left(\mathcal{A}^{\prime \prime}\right)=d_{H}\left(\mathcal{A}^{\prime \prime}, \mathcal{A}^{\prime}\right)+W\left(\mathcal{A}^{\prime}\right) \tag{4.1}
\end{equation*}
$$

Note that since $W$ is a solution of (2.5) and (2.6), then it is constant on each connected component of $\mathcal{A}_{\mathcal{H}}$. That is why hereafter we write $W(\mathcal{A}):=W(\xi), \xi \in \mathcal{A}$, for a connected component $\mathcal{A}$ of $\mathcal{A}_{\mathcal{H}}$. Since the distance function $d_{H}(x, y)$ satisfies the triangle inequality and $d_{H}\left(\mathcal{A}^{\prime \prime}, \mathcal{A}^{\prime}\right)+d_{H}\left(\mathcal{A}^{\prime}, \mathcal{A}^{\prime \prime}\right)>0$ for different components $\mathcal{A}^{\prime}, \mathcal{A}^{\prime \prime}$ of the Aubry set $\mathcal{A}_{H}$, (4.1) indeed defines a partial order relation.

Condition (2.11) assumes that there are finitely many different components of the Aubry set. It follows that there exists at least one minimal component $\mathcal{M}:=\mathcal{A}_{k_{0}}$ (such that, $\forall \mathcal{A}_{k} \neq \mathcal{M}$, either $\mathcal{M} \preceq \mathcal{A}_{k}$ or $\mathcal{M}$ and $\mathcal{A}_{k}$ are not comparable).

Now show that

$$
\begin{equation*}
W(x)=d_{H}(x, \mathcal{M})+W(\mathcal{M}) \text { in } U \cap \bar{\Omega}, \quad \text { where } U \text { is a neighbourhood of } \mathcal{M} . \tag{4.2}
\end{equation*}
$$

Indeed, otherwise, there is a sequence $x_{i} \rightarrow \mathcal{M}$ and a component $\mathcal{A}_{k} \neq \mathcal{M}$ such that $W\left(x_{i}\right)=d_{H}\left(x_{i}, \mathcal{A}_{k}\right)+W\left(\mathcal{A}_{k}\right)$. Then taking the limit, we derive $W(\mathcal{M})=d_{H}\left(\mathcal{M}, \mathcal{A}_{k}\right)+$ $W\left(\mathcal{A}_{k}\right)$, that is, $\mathcal{A}_{k} \preceq \mathcal{M}$ which contradicts the minimality of $\mathcal{M}$.

In what follows a component $\mathcal{M}$ such that (4.2) is satisfied will be called a significant component. We have shown that under condition (2.11) there is at least one significant component in the Aubry set $\mathcal{A}_{H}$.
Step II: Upper bound for eigenvalues. The crucial technical result in the proof of Theorem 2 is the following Lemma whose proof is presented in subsequent four Sections dealing separately with four possible cases of the structure of $\mathcal{M}$.

Lemma 5 Let $\mathcal{M}$ be a significant component of the Aubry set $\mathcal{A}_{H}$ satisfying either (2.12) or (2.13). Then for sufficiently small $\delta>0$, there are continuous functions $W_{\delta}^{ \pm}(x), W_{\delta, \varepsilon}^{ \pm}(x)$ and neighbourhoods $U_{\delta}$ of $\mathcal{M}$ such that

$$
W_{\delta}^{ \pm}(x)=0 \quad \text { on } \quad \mathcal{M}, \quad \text { and } \quad W_{\delta}^{-}(x)<W(x)-W(\mathcal{M})<W_{\delta}^{+}(x) \quad \text { in } \quad U_{\delta} \cap \bar{\Omega} \backslash \mathcal{M},
$$

$$
\begin{equation*}
W_{\delta, \varepsilon}^{ \pm} \in C^{2}\left(U_{\delta} \cap \bar{\Omega}\right), W_{\delta, \varepsilon}^{ \pm} \rightarrow W_{\delta}^{ \pm} \text {uniformly in } U_{\delta} \cap \bar{\Omega} \text { as } \varepsilon \rightarrow 0 \text {, and } \tag{4.3}
\end{equation*}
$$

$$
\begin{equation*}
\liminf _{\delta \rightarrow 0} \liminf _{\varepsilon \rightarrow 0, \xi_{\varepsilon} \rightarrow \mathcal{M}}\left(-a_{i j}\left(\xi_{\varepsilon}\right) \frac{\partial^{2} W_{\delta, \varepsilon}^{+}}{\partial x_{i} \partial x_{j}}\left(\xi_{\varepsilon}\right)+\frac{1}{\varepsilon} H\left(\nabla W_{\delta, \varepsilon}^{+}\left(\xi_{\varepsilon}\right), \xi_{\varepsilon}\right)+c\left(\xi_{\varepsilon}\right)\right) \geq \sigma(\mathcal{M}) \tag{4.4}
\end{equation*}
$$

$$
\begin{equation*}
\limsup _{\delta \rightarrow 0} \limsup _{\varepsilon \rightarrow 0, \xi_{\varepsilon} \rightarrow \mathcal{M}}\left(-a_{i j}\left(\xi_{\varepsilon}\right) \frac{\partial^{2} W_{\delta, \varepsilon}^{-}}{\partial x_{i} \partial x_{j}}\left(\xi_{\varepsilon}\right)+\frac{1}{\varepsilon} H\left(\nabla W_{\delta, \varepsilon}^{-}\left(\xi_{\varepsilon}\right), \xi_{\varepsilon}\right)+c\left(\xi_{\varepsilon}\right)\right) \leq \sigma(\mathcal{M}) \tag{4.5}
\end{equation*}
$$

Moreover, if $U_{\delta} \cap \partial \Omega \neq \emptyset$ then the functions $W_{\delta, \varepsilon}^{ \pm}$also satisfy $\frac{\partial W_{\delta, \varepsilon}^{+}}{\partial \nu}>0$ on $U_{\delta} \cap \partial \Omega$, and $\frac{\partial W_{\delta, \varepsilon}^{-}}{\partial v}<0$ on $U_{\delta} \cap \partial \Omega$.

Now, assuming that we know a minimal component $\mathcal{M}$ of the Aubry set $\mathcal{A}_{H}$, we can identify the limit $\lambda$ of eigenvalues $\lambda_{\varepsilon}$. Consider the difference $W_{\varepsilon}-W_{\delta, \varepsilon}^{-}$, where $W_{\delta, \varepsilon}^{-}$are test functions described in Lemma 5. By (4.3) the function $W-W_{\delta}^{-}-W(\mathcal{M})$ vanishes on $\mathcal{M}$ while it is strictly positive in a punctured neighbourhood of $\mathcal{M}$. Then, since $W_{\varepsilon}-W_{\delta, \varepsilon}^{-}$ converge uniformly to $W-W_{\delta}^{-}$as $\varepsilon \rightarrow 0$ in a neighbourhood of $\mathcal{M}$, there exists a sequence of local minima $\xi_{\varepsilon}$ of $W_{\varepsilon}-W_{\delta, \varepsilon}^{-}$such that $\xi_{\varepsilon} \rightarrow \mathcal{M}$. Moreover, if $\mathcal{M} \cap \partial \Omega \neq \emptyset$, then $\frac{\partial W_{\delta, \varepsilon}^{-}}{\partial \nu}<\frac{\partial W_{\varepsilon}}{\partial \nu}=0$ on $\partial \Omega$ (locally near $\mathcal{M}$ ) and therefore, $\xi_{\varepsilon} \in \Omega$ for sufficiently small $\varepsilon$. For such $\varepsilon$, we have

$$
\nabla W_{\varepsilon}=\nabla W_{\delta, \varepsilon}^{-} \quad \text { and } \quad-a_{i j} \frac{\partial^{2} W_{\varepsilon}}{\partial x_{i} \partial x_{j}} \leq-a_{i j} \frac{\partial^{2} W_{\delta, \varepsilon}^{-}}{\partial x_{i} \partial x_{j}} \quad \text { at } x=\xi_{\varepsilon}
$$

Therefore,

$$
\begin{aligned}
\lambda_{\varepsilon} & =-a_{i j}\left(\xi_{\varepsilon}\right) \frac{\partial^{2} W_{\varepsilon}}{\partial x_{i} \partial x_{j}}\left(\xi_{\varepsilon}\right)+\frac{1}{\varepsilon} H\left(\nabla W_{\varepsilon}\left(\xi_{\varepsilon}\right), \xi_{\varepsilon}\right)+c\left(\xi_{\varepsilon}\right) \\
& \leq-a_{i j}\left(\xi_{\varepsilon}\right) \frac{\partial^{2} W_{\delta, \varepsilon}^{-}}{\partial x_{i} \partial x_{j}}\left(\xi_{\varepsilon}\right)+\frac{1}{\varepsilon} H\left(\nabla W_{\delta, \varepsilon}^{-}\left(\xi_{\varepsilon}\right), \xi_{\varepsilon}\right)+c\left(\xi_{\varepsilon}\right) .
\end{aligned}
$$

Thus, we can use (4.5) here to pass first to the lim sup as $\varepsilon \rightarrow 0$ and then as $\delta \rightarrow 0$, this yields $\lim \sup _{\varepsilon \rightarrow 0} \lambda_{\varepsilon} \leq \sigma(\mathcal{M})$. Similarly one obtains the matching upper bound so that

$$
\begin{equation*}
\lim _{\varepsilon \rightarrow 0} \lambda_{\varepsilon}=\sigma(\mathcal{M}) \tag{4.6}
\end{equation*}
$$

However, since at this point $\mathcal{M}$ is unknown (it depends on $W$ and thus on the particular choice of a subsequence made in the beginning of the Section) equality (4.6) guarantees only the upper bound

$$
\begin{equation*}
\underset{\varepsilon \rightarrow 0}{\lim \sup } \lambda_{\varepsilon} \leq \max \left\{\sigma\left(\mathcal{A}_{k}\right) ; \mathcal{A}_{k} \subset \mathcal{A}_{H}\right\}, \tag{4.7}
\end{equation*}
$$

where the $\lim _{\sup _{\varepsilon \rightarrow 0}}$ is taken over the whole family $\left\{\lambda_{\varepsilon}, \varepsilon>0\right\}$.
Step III: Lower bound for eigenvalues. Consider a component $\mathcal{A}$ of the Aubry set $\mathcal{A}_{H}$ such that $\sigma(\mathcal{A})=\max \left\{\sigma\left(\mathcal{A}_{k}\right) ; \mathcal{A}_{k} \subset \mathcal{A}_{H}\right\}$. Introduce a smooth function $\rho(x)$ such that

$$
\rho(x) \geq 0 \text { in } \bar{\Omega}, \quad \rho(x)=0 \text { in a neighbourhood of } \mathcal{A} \text {, and } \rho(x)>0 \text {, when } x \in \mathcal{A}_{H} \backslash \mathcal{A}
$$

and consider the first eigenvalue $\bar{\lambda}_{\varepsilon}$ of the following auxiliary eigenvalue problem:

$$
\begin{equation*}
\varepsilon a_{i j}(x) \frac{\partial^{2} \bar{u}_{\varepsilon}}{\partial x_{i} \partial x_{j}}+b_{i}(x) \frac{\partial \bar{u}_{\varepsilon}}{\partial x_{i}}+\left(c(x)-\frac{1}{\varepsilon} \rho(x)\right) \bar{u}_{\varepsilon}=\bar{\lambda}_{\varepsilon} \bar{u}_{\varepsilon} \quad \text { in } \Omega, \tag{4.8}
\end{equation*}
$$

with the Neumann condition $\frac{\partial \bar{u}_{\varepsilon}}{\partial \nu}=0$ on $\partial \Omega$. By the Krein-Rutman theorem, the eigenvalue $\bar{\lambda}_{\varepsilon}$ is real and of multiplicity one, and $\bar{u}_{\varepsilon}$ being normalized by $\max _{\Omega} \bar{u}_{\varepsilon}=1$ satisfies $\bar{u}_{\varepsilon}>0$ in $\Omega$. Note that the adjoint problem also has a sign preserving eigenfunction.

Then it follows that

$$
\begin{equation*}
\bar{\lambda}_{\varepsilon} \leq \lambda_{\varepsilon} \tag{4.9}
\end{equation*}
$$

Indeed, otherwise, we have
$\varepsilon a_{i j}(x) \frac{\partial^{2} u_{\varepsilon}}{\partial x_{i} \partial x_{j}}+b_{i}(x) \frac{\partial u_{\varepsilon}}{\partial x_{i}}+\left(c(x)-\frac{1}{\varepsilon} \rho(x)\right) u_{\varepsilon}-\bar{\lambda}_{\varepsilon} u_{\varepsilon}=-\left(\bar{\lambda}_{\varepsilon}-\lambda_{\varepsilon}+\frac{1}{\varepsilon} \rho(x)\right) u_{\varepsilon}<0 \quad$ in $\Omega$.
On the other hand, by Fredholm's theorem the right-hand side in (4.10) must be orthogonal (in $L^{2}(\Omega)$ ) to any eigenfunction of the problem adjoint to (4.8). Since the latter problem has a sign preserving eigenfunction, we arrive at a contradiction which proves (4.9).

Let $\bar{W}_{\varepsilon}:=-\varepsilon \log \bar{u}_{\varepsilon}$ be the scaled logarithmic transform of $\bar{u}_{\varepsilon}$, i.e. $\bar{u}_{\varepsilon}=e^{-\bar{W}_{\varepsilon} / \varepsilon}$. Following the line of Section 3 one can show that, up to extracting a subsequence, functions $\bar{W}_{\varepsilon}$ converge (uniformly in $\bar{\Omega}$ ) to a viscosity solution $\bar{W}$ of the problem

$$
\begin{equation*}
H(\nabla \bar{W}(x), x)-\rho(x)=\Lambda \quad \text { in } \Omega \tag{4.11}
\end{equation*}
$$

with the boundary condition $\frac{\partial \bar{W}}{\partial v}=0$, where $\Lambda=\lim _{\varepsilon \rightarrow 0} \varepsilon \bar{\lambda}_{\varepsilon}$. Note that the argument in Lemma 3 yields now bounds of the form $-\frac{C}{\varepsilon} \leq \bar{\lambda}_{\varepsilon} \leq C$ with some $C>0$ independent
of $\varepsilon$. Nevertheless, these bounds are sufficient to derive problem (4.11) with the Neumann boundary condition.

Moreover, since $\rho=0$ in a neighbourhood of $\mathcal{A}$ one can show that $\Lambda=0$ using testing curves $\eta$ from (2.10) in the variational representation for the additive eigenvalue $\Lambda$ (see [14]),

$$
\Lambda=-\lim _{T \rightarrow \infty} \inf \left\{\frac{1}{T} \int_{0}^{T}(L(-v, \eta)+\rho(\eta)) \mathrm{d} t ; \eta \text { solves (2.8) with } \eta(0)=x \in \bar{\Omega}\right\}
$$

This implies, in particular, that

$$
\bar{W}(x)=d_{H}(x, \mathcal{A}) \text { in a neighbourhood of } \mathcal{A},
$$

where $d_{H}(x, y)$ is the distance function given by (2.9). Then arguing as in second step we obtain

$$
\bar{\lambda}_{\varepsilon} \rightarrow \sigma(\mathcal{A})
$$

Thanks to (4.9) this yields the lower bound $\lim \inf \lambda_{\varepsilon} \geq \max \left\{\sigma\left(\mathcal{A}_{k}\right) ; \mathcal{A}_{k} \in \mathcal{A}_{H}\right\}$ complementary to (4.7). Thus formula (2.18) is proved.
Step IV: Selection of the solution of (2.5)-(2.6). Let us assume now that the maximum in (2.18) is attained at exactly one component $\mathcal{M}$. Then comparing (2.18) with (4.6) we see that $\mathcal{M}$ is the unique significant component in $\mathcal{A}_{H}$; therefore, it is the only minimal component of $\mathcal{A}_{H}$ with respect to the order relation $\preceq$. Thus $\mathcal{M}$ is the least component of $\mathcal{A}_{H}$. It follows that $W\left(\mathcal{A}_{k}\right)-W(\mathcal{M})=d_{H}\left(\mathcal{A}_{k}, \mathcal{M}\right)$ for every $\mathcal{A}_{k} \subset \mathcal{A}_{H}$. Then by (2.7), the representation $W(x)=d_{H}(x, \mathcal{M})+W(\mathcal{M})$ holds in $\bar{\Omega}$. Finally, since $\min _{\bar{\Omega}} W(x)=$ $\lim _{\varepsilon \rightarrow 0} \min _{\bar{\Omega}} W_{\varepsilon}(x)=0$, we have $W(\mathcal{M})=0$, i.e. $W(x)=d_{H}(x, \mathcal{M})$. Theorem 2 is proved.

## 5. Construction of test functions: case of fixed points in $\boldsymbol{\Omega}$

The central part in the proof of Theorem 2 is the construction of test functions satisfying the conditions of Lemma 5 for different types of components of the Aubry set $\mathcal{A}_{H}$. Consider first the case when a fixed point $\xi \in \Omega$ of the $\mathrm{ODE} \dot{x}=b(x)$ is a significant component of $\mathcal{A}_{H}$. We can assume that $W(\xi)=0$, subtracting an appropriate constant if necessary. Then $W(x)$ is given by

$$
\begin{equation*}
W(x)=d_{H}(x, \xi) \quad \text { in a neighbourhood } U(\xi) \text { of } \xi . \tag{5.1}
\end{equation*}
$$

We begin by studying the local behaviour of $W(x)$ near $\xi$. Consider for sufficiently small $z$ the ansatz

$$
\begin{equation*}
W(z+\xi)=\Gamma_{i j} z_{i} z_{j}+o\left(|z|^{2}\right) \tag{5.2}
\end{equation*}
$$

with a symmetric $N \times N$ matrix $\Gamma$. After substituting (5.2) into (2.5) and collecting quadratic terms in the resulting relation, we are led to the Riccati matrix equation

$$
\begin{equation*}
4 \Gamma Q \Gamma-\Gamma B-B^{*} \Gamma=0, \tag{5.3}
\end{equation*}
$$

where $Q=\left(a_{i j}(\xi)\right)_{i, j=\overline{1, N}}, B=\left(\frac{\partial b_{i}}{\partial x_{j}}(\xi)\right)_{i, j=\overline{1, N}}$.
Next, we show that (5.2) holds with $\Gamma$ being the maximal symmetric solution of (5.3); for existence of such a solution see, e.g. [21] or [22] . To this end,
consider the solution $D$ of the Lyapunov matrix equation

$$
\begin{equation*}
D(4 \Gamma Q-B)+(4 \Gamma Q-B)^{*} D=2 I \tag{5.4}
\end{equation*}
$$

given by

$$
\begin{equation*}
D=2 \int_{-\infty}^{0} e^{(4 \Gamma Q-B)^{*} t} e^{(4 Q \Gamma-B) t} \mathrm{~d} t \tag{5.5}
\end{equation*}
$$

By Theorem 9.1 .3 in [21], all the eigenvalues of the matrix $4 Q \Gamma-B$ have positive real parts, so that the integral in (5.5) does converge. Set

$$
\Gamma_{\delta}^{ \pm}=\Gamma \pm \delta D
$$

Then $\Gamma_{\delta}^{-}$satisfies

$$
\begin{equation*}
4 \Gamma_{\delta}^{-} Q \Gamma_{\delta}^{-}-\Gamma_{\delta}^{-} B-B^{*} \Gamma_{\delta}^{-} \leq-\delta I \tag{5.6}
\end{equation*}
$$

for sufficiently small $\delta>0$.
Introduce the quadratic function $W_{\delta}^{-}(x):=\Gamma_{\delta}^{-}(x-\xi) \cdot(x-\xi)$. Thanks to (5.6) this function satisfies

$$
\begin{equation*}
H\left(\nabla W_{\delta}^{-}(x), x\right) \leq-\frac{\delta}{2}|x-\xi|^{2} \quad \text { in a neighbourhood of } \xi \tag{5.7}
\end{equation*}
$$

This yields the following result whose proof is identical to the proof of Lemma 16 in [12] (see also the arguments in the proof of Lemma 8 below).

Lemma 6 The strict pointwise inequality $W_{\delta}^{-}(x)<W(x)$ holds in a punctured neighbourhood of $\xi$ for sufficiently small $\delta>0$.

Next consider the function $W_{\delta}^{+}(x):=\Gamma_{\delta}^{+}(x-\xi) \cdot(x-\xi)$.
Lemma 7 The strict pointwise inequality $W_{\delta}^{+}(x)>W(x)$ holds in a punctured neighbourhood of $\xi$ for sufficiently small $\delta>0$.

Proof According to (5.1), the following inequality holds

$$
W(x) \leq \int_{0}^{t} L(-v(\tau), \xi+\eta(\tau)) \mathrm{d} \tau
$$

for every control $v(\tau)$ such that the solution of the ODE

$$
\dot{\eta}(\tau)=v(\tau), \quad \eta(0)=z:=x-\xi
$$

vanishes at the final time $t$ and remains in a small neighbourhood of 0 for any $0 \leq \tau \leq t$. We can take the final time $t=+\infty$ and construct $v(\tau)$ by setting $v(\tau)=-(4 Q \Gamma-B) \eta(\tau)$, where $\eta(\tau)$ in the solution of the ODE

$$
\dot{\eta}=-(4 Q \Gamma-B) \eta, \quad \eta(0)=z
$$

As already mentioned (see Theorem 9.1.3. in [21]), all the eigenvalues of the matrix $4 Q \Gamma-B$ have positive real parts, therefore $|\eta(\tau)| \leq C|z|$ and $\eta(\tau) \rightarrow 0$ as $\tau \rightarrow+\infty$. Moreover, the latter convergence is exponentially fast.

Thus we have

$$
\begin{aligned}
L(-v(\tau), \eta(\tau)+\xi) & =\frac{1}{4} a^{i j}(\xi+\eta)\left(-\dot{\eta}_{i}+b_{i}(\eta)\right)\left(-\dot{\eta}_{i}+b_{i}(\eta)\right) \\
& =\frac{1}{4} a^{i j}(\xi)\left(-\dot{\eta}_{i}+B_{i k} \eta_{k}\right)\left(-\dot{\eta}_{j}+B_{j l} \eta_{l}\right)+O\left(|\eta|^{3}\right),
\end{aligned}
$$

where $\left(a^{i j}(x)\right)_{i, j=\overline{1, N}}$ denotes the matrix inverse to $\left(a_{i j}(x)\right)_{i, j=\overline{1, N}}$. Next, recall that $Q_{i j}=$ $a_{i j}(\xi)$ and that $\Gamma$ solves (5.3). Taking this into account we obtain

$$
\begin{aligned}
\int_{0}^{\infty} L(\eta+\xi,-v(\tau)) & =\frac{1}{4} \int_{0}^{\infty} a^{i j}(\xi)\left(-\dot{\eta}_{i}+B_{i k} \eta_{k}\right)\left(-\dot{\eta}_{j}+B_{j l} \eta_{l}\right)+O\left(|z|^{3}\right) \\
& =-2 \int_{0}^{\infty} \Gamma \eta \cdot \dot{\eta} \mathrm{d} \tau+\int_{0}^{\infty} \Gamma \eta \cdot(\dot{\eta}+B \eta) \mathrm{d} \tau+O\left(|z|^{3}\right) \\
& =\Gamma z \cdot z+\int_{0}^{\infty} \eta \cdot\left(-4 \Gamma Q \Gamma+\Gamma B+B^{*} \Gamma\right) \eta \mathrm{d} \tau+O\left(|z|^{3}\right) \\
& =\Gamma_{i j} z_{i} z_{j}+O\left(|z|^{3}\right) .
\end{aligned}
$$

Finally, since $\Gamma_{\delta}^{+}=\Gamma+\delta D$ with $D>0$, then for sufficiently small $z \neq 0$ we have

$$
W(z+\xi) \leq \Gamma z \cdot z+O\left(|z|^{3}\right)<\Gamma_{\delta}^{+} z \cdot z .
$$

Lemmas 6 and 7 show that functions $W_{\delta}^{ \pm}$do satisfy conditions of Lemma 5. To complete the proof of Lemma 5 in the case of $\mathcal{M}$ being a fixed point in $\Omega$, we define functions $W_{\delta, \varepsilon}^{ \pm}$ simply by setting $W_{\delta, \varepsilon}^{ \pm}:=W_{\delta}^{ \pm}$. Thanks to (5.7) we have
$\limsup _{\varepsilon \rightarrow 0}\left(a_{i j}\left(\xi_{\varepsilon}\right) \frac{\partial^{2} W_{\delta, \varepsilon}^{-}}{\partial x_{i} \partial x_{j}}\left(\xi_{\varepsilon}\right)+\frac{1}{\varepsilon} H\left(\nabla W_{\delta, \varepsilon}^{-}\left(\xi_{\varepsilon}\right), \xi_{\varepsilon}\right)+c\left(\xi_{\varepsilon}\right)\right) \leq-2 a_{i j}(\xi)\left(\Gamma_{i j}-\delta D_{i j}\right)+c(\xi)$,
as soon as $\xi_{\varepsilon} \rightarrow \xi$ when $\varepsilon \rightarrow 0$. According to Proposition 20 in [12], $-2 a_{i j}(\xi) \Gamma_{i j}+c(\xi)=$ $\sigma(\{\xi\})$, thus (5.8) yields (4.5). Similarly one verifies that $W_{\delta, \varepsilon}^{+}$satisfies (4.4).

## 6. Construction of test functions: case of fixed points on $\partial \boldsymbol{\Omega}$

Consider now the case of significant component of the Aubry set $\mathcal{A}_{H}$ being a hyperbolic fixed point $\xi$ of the $\mathrm{ODE} \dot{x}=b_{\tau}(x)$ on $\partial \Omega$, where $b_{\tau}(x)$ denotes the tangential component of the vector field $b(x)$ on $\partial \Omega$. As above, without loss of generality, we assume that $W(\xi)=0$.

It is convenient to introduce local coordinates near $\partial \Omega$ so that $x=X\left(z_{1}, \ldots, z_{N}\right)$ with $z_{N}=z_{N}(x)$ being the distance from $x$ to $\partial \Omega\left(z_{N}(x)>0\right.$ if $\left.x \in \Omega\right)$ and $z^{\prime}=\left(z_{1}, \ldots, z_{N-1}\right)$ representing coordinates on $\partial \Omega$ in a neighbourhood of the point $\xi$. The latter coordinates are chosen so that the map $X\left(z^{\prime}, z_{N}\right)$ is $C^{2}$-smooth and $z^{\prime}(\xi)=0$. Moreover, the matrix $\left(\frac{\partial X_{i}}{\partial z_{j}}\right)_{i, j=\overline{1, N}}$ is orthogonal when $z^{\prime}=0$ and $z_{N}=0$ (at the point $\xi$ ). In these new variables, Equations (2.5) and (2.2) read

$$
\begin{equation*}
S\left(\nabla_{z} W, z\right)=0 \tag{6.1}
\end{equation*}
$$

and

$$
\begin{equation*}
-\varepsilon a_{i j}(X(z)) \mathcal{T}_{k i}^{-1}(z) \frac{\partial}{\partial z_{k}}\left(\mathcal{T}_{l j}^{-1}(z) \frac{\partial W_{\varepsilon}}{\partial z_{l}}\right)+S\left(\nabla_{z} W_{\varepsilon}, z\right)=\varepsilon\left(\lambda_{\varepsilon}-c(X(z))\right), \tag{6.2}
\end{equation*}
$$

where

$$
S(p, z)=a_{i j}(X(z)) \mathcal{T}_{k i}^{-1}(z) \mathcal{T}_{l j}^{-1}(z) p_{k} p_{l}-b_{i}(X(z)) \mathcal{T}_{k i}^{-1}(z) p_{k}
$$

and $\left(\mathcal{T}_{i j}^{-1}(z)\right)_{i, j=\overline{1, N}}$ is the inverse matrix to $\left(\frac{\partial X_{i}}{\partial z_{j}}(z)\right)_{i, j=\overline{1, N}}$. Note that according to hypothesis (2.13)

$$
\begin{equation*}
b_{i}(X(z)) \mathcal{T}_{N i}^{-1}(z)<0 \text { for sufficiently small }|z| . \tag{6.3}
\end{equation*}
$$

Like in Section 5 we construct the leading term of the asymptotic expansion of $W$ near the fixed point $\xi$ in the form of a quadratic function. Taking into account the boundary condition $\frac{\partial W}{\partial z_{N}}=0$ (that is (2.6) rewritten in aforementioned local coordinates), we write down the following ansatz

$$
W\left(X\left(z^{\prime}, z_{N}\right)\right)=\tilde{\Gamma}_{i j} z_{i}^{\prime} z_{j}^{\prime}+o\left(|z|^{2}+z_{N}^{2}\right) .
$$

with a symmetric $(N-1) \times(N-1)$ matrix $\widetilde{\Gamma}$ satisfying the Riccati equation

$$
\begin{equation*}
4 \tilde{\Gamma} \tilde{Q} \tilde{\Gamma}-\tilde{\Gamma} \tilde{B}-\tilde{B}^{*} \tilde{\Gamma}=0, \tag{6.4}
\end{equation*}
$$

where $\tilde{Q}=\left(a_{i j}(\xi) \mathcal{T}_{k i}^{-1}(0) \mathcal{T}_{l j}^{-1}(0)\right)_{k, l=\overline{1, N-1}}$ and $\tilde{B}=\left(\mathcal{T}_{k i}^{-1}(0) \frac{\partial b_{i}}{\partial x_{j}}(\xi) \frac{\partial X_{j}}{\partial z l}(0)\right)_{k, l=\overline{1, N-1}}$. Note that $\tilde{B}$ is nothing but the matrix in the $\mathrm{ODE} \dot{z}^{\prime}=\tilde{B} z^{\prime}$ obtained by linearizing the ODE $\dot{x}=b_{\tau}(x)$ near $\xi$ in the local coordinates $z^{\prime}=\left(z_{1}^{\prime}, \ldots, z_{N-1}^{\prime}\right)$ on $\partial \Omega$.

Let $\tilde{\Gamma}$ be the maximal symmetric solution of (6.4), and let $\tilde{D}$ be a solution of the Lyapunov matrix equation

$$
\begin{equation*}
\tilde{D}(4 \tilde{\Gamma} \tilde{Q}-\tilde{B})+(4 \tilde{\Gamma} \tilde{Q}-\tilde{B})^{*} \tilde{D}=2 I . \tag{6.5}
\end{equation*}
$$

By Theorem 9.1.3 in [21], all the eigenvalues of the matrix $4 \tilde{\Gamma} \tilde{Q}-\tilde{B}$ have positive real parts, therefore (6.5) has the unique solution $\tilde{D}$ given by

$$
\tilde{D}=2 \int_{-\infty}^{0} e^{(4 \tilde{\Gamma} \tilde{Q}-\tilde{B})^{*} t} e^{(4 \tilde{\Gamma} \tilde{Q}-\tilde{B}) t} \mathrm{~d} t
$$

which is a symmetric positive definite matrix. Now introduce functions

$$
\begin{equation*}
W_{\delta}^{ \pm}\left(z^{\prime}, z_{N}\right)=(\tilde{\Gamma} \pm \delta \tilde{D})_{i j} z_{i}^{\prime} z_{j}^{\prime} \pm \delta z_{N}^{2} \tag{6.6}
\end{equation*}
$$

depending on the parameter $\delta>0$.
Lemma 8 Let $\delta>0$ be sufficiently small. Then, for small $|z| \neq 0$ such that $X(z) \in \bar{\Omega}$, we have

$$
\begin{equation*}
W_{\delta}^{-}(z)<W(X(z))<W_{\delta}^{+}(z) . \tag{6.7}
\end{equation*}
$$

Proof By virtue of the definition of $W_{\delta}^{ \pm}$, it suffices to prove (6.7) with non-strict inequalities in place of strict ones and then pass to slightly bigger $\delta$.

The proof of the inequality $W_{\delta}^{-} \leq W$ is based on the following two facts. First, we use the fact that $W(x)=d_{H}(x, \xi)$ in a neighbourhood of $\xi$. Moreover, for a given $\delta^{\prime}>0$, there
exists $\delta>0$ such that if $|x-\xi|<\delta$, then the minimization in (2.9) is actually restricted to testing curves $\eta(\tau)$ which do not leave the set $\left\{|\eta-\xi|<\delta^{\prime}\right\}$ (otherwise, arguing as in [12, Lemma 19], one can show that $\xi$ is not an isolated point of the Aubry set $\mathcal{A}_{H}$, contradicting (2.11)). Second, considering, with a little abuse of notation, $W_{\delta}^{-}(x)=W_{\delta}^{-}\left(X^{-1}(x)\right)$, we have for sufficiently small $\delta>0$

$$
\begin{equation*}
H\left(\nabla W_{\delta}^{-}, x\right) \leq-\delta|x-\xi|^{2} \text { in } \Omega, \text { and } \frac{\partial W_{\delta}^{-}}{\partial v}=0 \text { on } \partial \Omega \tag{6.8}
\end{equation*}
$$

when $|x-\xi|<\delta^{\prime}$ with some $\delta^{\prime}>0$ independent of $\delta$. This follows from the construction (6.6) of $W_{\delta}^{ \pm}$and (6.4), (6.5), also taking into account (6.3).

Assume that $|x-\xi|<\delta$, and let $\eta(\tau)$ be a solution of (2.8) satisfying $\eta(0)=x, \eta(t)=\xi$ with a control $v(\tau)$ such that $|\eta(\tau)-\xi|<\delta^{\prime}$ for all $0 \leq \tau \leq t$. Then

$$
W_{\delta}^{-}(x)=-\int_{0}^{t} \nabla W_{\delta}^{-}(\eta) \cdot \dot{\eta} \mathrm{d} \tau=\int_{0}^{t} \nabla W_{\delta}^{-}(\eta) \cdot(-v(\tau)) \mathrm{d} \tau,
$$

where we have used the fact that $\frac{\partial W_{\delta}^{-}}{\partial \nu}=0$ on $\partial \Omega$. It follows by Fenchel's inequality $p \cdot(-v) \leq L(-v, \eta)+H(p, \eta)$ that

$$
W_{\delta}^{-}(x) \leq \int_{0}^{t} L(-v, \eta) \mathrm{d} \tau+\int_{0}^{t} H\left(\nabla W_{\delta}^{-}, \eta\right) \mathrm{d} \tau \leq \int_{0}^{t} L(-v, \eta) \mathrm{d} \tau .
$$

Therefore by (2.9), we obtain $W_{\delta}^{-}(x) \leq W(x)$.
In order to prove the second inequality in (6.7) for a given $x=X\left(z^{\prime}, z_{N}\right)$, we construct a test curve $\eta(\tau)$ first on a small interval $(0, \Delta t)$ by setting $\eta(\tau)=X\left(z^{\prime}, \zeta_{N}(\tau)\right), \zeta_{N}(\tau)$ being the solution of $\operatorname{ODE} \dot{\zeta}_{N}(\tau)=b_{i}\left(X\left(z^{\prime}, \zeta_{N}\right)\right) \mathcal{T}_{N i}^{-1}\left(z^{\prime}, \zeta_{N}\right)$ with the initial condition $\zeta_{N}(0)=z_{N}$, and choosing $\Delta t$ from the conditions $\zeta_{N}(\Delta t)=0, \zeta_{N}(\tau)>0$ for $\tau<\Delta t$. Thanks to (6.3) we have $\Delta t=O\left(z_{N}\right)$. Then, since

$$
\begin{aligned}
\dot{\eta}_{i} & =\frac{\partial X_{i}}{\partial z_{N}}\left(z^{\prime}, \zeta_{N}\right) b_{k}(\eta) \mathcal{T}_{N k}^{-1}\left(z^{\prime}, \zeta_{N}\right) \\
& =\frac{\partial X_{i}}{\partial z_{j}}\left(z^{\prime}, \zeta_{N}\right) b_{k}(\eta) \mathcal{T}_{j k}^{-1}\left(z^{\prime}, \zeta_{N}\right)-\frac{\partial X_{i}}{\partial z_{j}^{\prime}}\left(z^{\prime}, \zeta_{N}\right) b_{k}(\eta) \mathcal{T}_{j k}^{-1}\left(z^{\prime}, \zeta_{N}\right)=b_{i}(\eta)+O(|z|)
\end{aligned}
$$

(recall that the tangential component $b_{\tau}$ on $\partial \Omega$ vanishes at the point $\xi$ ), we obtain

$$
\begin{equation*}
\int_{0}^{\Delta t} L(-\dot{\eta}, \eta) \mathrm{d} \tau=O\left(|z|^{3}\right) \tag{6.9}
\end{equation*}
$$

Next, we construct $\eta(\tau)$ for $\tau>\Delta t$ which connects point $X\left(z^{\prime}, 0\right)$ to $\xi$. Following closely the line of Lemma 7 we introduce $\zeta^{\prime}(\tau)$ by solving the equation $\dot{\zeta}^{\prime}(\tau)=-(4 \tilde{Q} \tilde{\Gamma}-\tilde{B}) \zeta^{\prime}$ with the initial condition $\zeta^{\prime}(\Delta t)=z^{\prime}$ and set $\eta(\tau)=X\left(\zeta^{\prime}(\tau), 0\right)$. Then $\eta(\tau)$ solves (2.8) for $\tau>\Delta t$ with

$$
v(\tau):=\dot{\eta}(\tau)+v(\eta) b_{v}(\eta)-4 \frac{\partial X}{\partial z_{N}}\left(\zeta^{\prime}, 0\right) \mathcal{T}_{N i}^{-1}(0) a_{i j}(\xi) \mathcal{T}_{l j}^{-1}(0) \tilde{\Gamma}_{l m} \zeta_{m}^{\prime}
$$

(note that $\frac{\partial X}{\partial z_{N}}\left(\zeta^{\prime}, 0\right)=-\nu(\eta)$ and $b_{v}(\eta)+4 \mathcal{T}_{N i}^{-1}(0) a_{i j}(\xi) \mathcal{T}_{l j}^{-1}(0) \Gamma_{l m} \zeta_{m}^{\prime}>0$ as soon as $\left|z^{\prime}\right|$ is sufficiently small) and using (6.4), we obtain

$$
\begin{align*}
\int_{\Delta t}^{\infty} L(- & v(\tau), \eta(\tau)) \mathrm{d} \tau=\frac{1}{4} \int_{\Delta t}^{\infty} a^{i j}(\xi)\left(-v_{i}+b_{i}\right)\left(-v_{j}+b_{j}\right) \mathrm{d} \tau+O\left(|z|^{3}\right) \\
= & 4 \int_{\Delta t}^{\infty} a^{i j}(\xi)\left(\frac{\partial X_{i}}{\partial z_{k}^{\prime}}(0) \tilde{Q}_{k l} \tilde{\Gamma}_{l m} \zeta_{m}^{\prime}+\frac{\partial X_{i}}{\partial z_{N}}(0) \mathcal{T}_{N k}^{-1}(0) a_{k l}(\xi) \mathcal{T}_{m l}^{-1}(0) \tilde{\Gamma}_{m n} \zeta_{n}^{\prime}\right) \\
& \times\left(\frac{\partial X_{j}}{\partial z_{k}^{\prime}}(0) \tilde{Q}_{k l} \tilde{\Gamma}_{l m} \zeta_{m}^{\prime}+\frac{\partial X_{j}}{\partial z_{N}}(0) \mathcal{T}_{N k}^{-1}(0) a_{k l}(\xi) \mathcal{T}_{m l}^{-1}(0) \tilde{\Gamma}_{m n} \zeta_{n}^{\prime}\right) \mathrm{d} \tau+O\left(|z|^{3}\right) \\
= & 4 \int_{\Delta t}^{\infty} \tilde{\Gamma} \zeta^{\prime} \cdot \tilde{Q} \tilde{\Gamma} \zeta^{\prime} \mathrm{d} \tau+O\left(|z|^{3}\right) \\
= & -2 \int_{\Delta t}^{\infty} \tilde{\Gamma} \zeta^{\prime} \cdot \dot{\zeta}^{\prime} \mathrm{d} \tau+\int_{\Delta t}^{\infty} \tilde{\Gamma} \zeta^{\prime} \cdot\left(\dot{\zeta}^{\prime}+\tilde{B} \zeta^{\prime}\right) \mathrm{d} \tau+O\left(|z|^{3}\right)=\tilde{\Gamma}_{i j} z_{i}^{\prime} z_{j}^{\prime}+O\left(|z|^{3}\right) \tag{6.10}
\end{align*}
$$

The required upper bound $W \leq W_{\delta}^{+}$now follows from (6.9) to (6.10).
Thus functions $W_{\delta}^{ \pm}$satisfy conditions of Lemma 5; moreover, it follows from (6.6) in conjunction with (6.4), (6.5), taking also into account (6.3), that

$$
S\left(\nabla_{z} W_{\delta}^{+}(z), z\right) \geq 0 \quad \text { and } \quad S\left(\nabla_{z} W_{\delta}^{-}(z), z\right) \leq 0 \quad \text { when }|z| \text { is sufficiently small. }
$$

Then we set

$$
W_{\delta, \varepsilon}^{ \pm}\left(z^{\prime}, z_{N}\right)=W_{\delta}^{ \pm}\left(z^{\prime}, z_{N}\right) \mp \varepsilon^{2} z_{N}
$$

and verify (similarly to the case of interior fixed points) that conditions (4.4) and (4.5) are satisfied. Additionally, we have $\mp \frac{\partial W_{\delta, \varepsilon}^{ \pm}}{\partial z_{N}}\left(z^{\prime}, 0\right)>0$, i.e. $\pm \frac{\partial W_{\delta, \varepsilon}^{ \pm}}{\partial \nu}>0$ on $\partial \Omega$.

## 7. Construction of test functions: case of limit cycles in $\Omega$

We proceed with the case when a significant component of the Aubry set $\mathcal{A}_{H}$ is a limit cycle, assuming first that it is situated entirely inside $\Omega$. Namely, let $\xi(t)$ be a periodic solution of the ODE $\dot{\xi}=b(\xi)$ whose minimal period is $P>0$. We assume that $\mathcal{C}=\{\xi(t): t \in$ $[0, P)\} \subset \Omega, b(x) \neq 0$ on $\mathcal{C}$ and $\mathcal{C}$ is a hyperbolic limit cycle, i.e. the linearized Poincaré map associated to this cycle has no eigenvalues on the unit circle. In order to study the local behaviour of
$W$ near the cycle $\mathcal{C}$, perform a $C^{2}$-smooth change of coordinates $x=X\left(z_{1}, \ldots, z_{N-1}\right.$, $\left.z_{N}\right)$ with $z_{N}$ representing the arc length along the cycle and $z^{\prime}=\left(z_{1}, \ldots, z_{N-1}\right)$ being some fixed Cartesian coordinates in the hyperplanes orthogonal to the cycle. Also, we assume that $\mathcal{C}$ is oriented by the tangent vector $b(\xi) /|b(\xi)|$, and $z^{\prime}=0$ on $\mathcal{C}$. With this change of coordinates, Equations (2.5) and (2.2) take the form similar to (6.1) and (6.2).

Assuming as above that $W(\mathcal{C})=0$, we postulate in the vicinity of the cycle (for sufficiently small $\left|z^{\prime}\right|$ ) the following ansatz for $W$ :

$$
\begin{equation*}
W\left(X\left(z^{\prime}, z_{N}\right)\right)=\bar{\Gamma}_{i j}(t) z_{i}^{\prime} z_{j}^{\prime}+o\left(\left|z^{\prime}\right|^{2}\right), \tag{7.1}
\end{equation*}
$$

where $t$ refers to the parametrization of the cycle determined by the equation $\dot{\xi}=b(\xi)$, $t \in[0, P)$.

Substitute $W$ in (6.1) to find, after collecting quadratic terms and neglecting higher order terms,

$$
\begin{equation*}
\dot{\bar{\Gamma}}=4 \bar{\Gamma} \bar{Q} \bar{\Gamma}-\bar{\Gamma} \bar{B}-\bar{B}^{*} \bar{\Gamma} \tag{7.2}
\end{equation*}
$$

where $\bar{Q}(t)$ and $\bar{B}(t)$ are $P$-periodic $(N-1) \times(N-1)$ matrices whose entries are given by

$$
\begin{align*}
\bar{Q}_{k l}(t)= & a_{i j}(\xi(t)) \mathcal{T}_{k i}^{-1}\left(0, z_{N}(\xi(t))\right) \mathcal{T}_{l j}^{-1}\left(0, z_{N}(\xi(t))\right)  \tag{7.3}\\
\bar{B}_{k l}(t)= & \mathcal{T}_{k i}^{-1}\left(0, z_{N}(\xi(t))\right) \frac{\partial b_{i}}{\partial x_{j}}(\xi(t)) \frac{\partial X_{j}}{\partial z_{l}}\left(0, z_{N}(\xi(t))\right) \\
& -\mathcal{T}_{k j}^{-1}\left(0, z_{N}(\xi(t))\right) \frac{\mathrm{d}}{\mathrm{~d} t}\left(\frac{\partial X_{j}}{\partial z_{l}}\left(0, z_{N}(\xi(t))\right)\right) \tag{7.4}
\end{align*}
$$

Recall that $\mathcal{T}_{i j}^{-1}\left(z^{\prime}, z_{N}\right)$ denote the entries of the matrix inverse to $\left(\frac{\partial X_{i}}{\partial z_{j}}\left(z^{\prime}, z_{N}\right)\right)_{i, j=\overline{1, N}}$ and for brevity, abusing slightly the notation, we set

$$
\begin{equation*}
\mathcal{T}_{i j}^{-1}(t):=\mathcal{T}_{i j}^{-1}\left(0, z_{N}(\xi(t)), \quad \mathcal{T}_{i j}(t)=\frac{\partial X_{i}}{\partial z_{j}}\left(0, z_{N}(\xi(t))\right)\right. \tag{7.5}
\end{equation*}
$$

The matrix $\bar{Q}(t)$ being positive definite, it is known [22] that Riccati Equation (7.2) has a maximal symmetric P-periodic solution $\bar{\Gamma}(t)$. We next show that (7.1) does hold with the mentioned maximal solution $\bar{\Gamma}(t)$ under our standing hyperbolicity assumption on $\mathcal{C}$. Note that the $\mathrm{ODE} \dot{z}^{\prime}=\bar{B} z^{\prime}$ corresponds to the linearization of $\dot{x}=b(x)$ on the cycle $\mathcal{C}$ written in local coordinates; thus assuming the hyperbolicity of $\mathcal{C}$, we require that the fundamental solution of the ODE

$$
\frac{\partial \bar{\Phi}}{\partial t}(t, \tau)=\bar{B}(t) \bar{\Phi}(t, \tau), \quad \bar{\Phi}(\tau, \tau)=I
$$

evaluated at $t=\tau+P$, has no eigenvalues with absolute value equal to 1 .
Lemma 9 The following bound holds uniformly in $t \in[0, P)$ for sufficiently small $\left|z^{\prime}\right|$,

$$
\begin{equation*}
W\left(X\left(z^{\prime}, z_{N}(\xi(t))\right) \leq \bar{\Gamma}_{i j}(t) z_{i}^{\prime} z_{j}^{\prime}+C\left|z^{\prime}\right|^{3} \log \frac{1}{\left|z^{\prime}\right|}\right. \tag{7.6}
\end{equation*}
$$

Proof We make use of variational representation (2.9). A natural guess about optimal test curve in (2.9) is that its first $N-1$ local coordinates are given by (cf. Sections 5 and 6)

$$
\begin{equation*}
\dot{\zeta}^{\prime}(\tau)=(\bar{B}-4 \bar{Q} \bar{\Gamma}) \zeta^{\prime}(\tau), \quad \tau>t, \quad \zeta^{\prime}(t)=z^{\prime} \tag{7.7}
\end{equation*}
$$

Thanks to Theorem 5.4.15 in [22] the solutions of (7.7) are exponentially stable, i. e. $\left|\zeta^{\prime}(\tau)\right| \leq C e^{-\delta \tau}\left|z^{\prime}\right|$ for some $\delta>0$. The choice of the last local coordinate is a little bit involved. We set $\eta(\tau):=X\left(\zeta^{\prime}(\tau), z_{N}(\xi(\tau))+\zeta_{N}(\tau)\right)$ and want to choose $\zeta_{N}(\tau)$ in such a way that $\left|\zeta_{N}(\tau)\right|<C\left|z^{\prime}\right|$, and

$$
\begin{equation*}
b_{i}(\eta(\tau))-\dot{\eta}_{i}(\tau)=4 a_{i j}(\xi(\tau)) \mathcal{T}_{l j}^{-1}(0, \xi(\tau)) \bar{\Gamma}_{l m}(\tau) \zeta_{m}^{\prime}(\tau)+O\left(\left|z^{\prime}\right|^{2}\right) \tag{7.8}
\end{equation*}
$$

We skip for a moment the proof of the existence of such $\zeta_{N}$. It will be given later on.

Considering (7.8) we obtain

$$
\begin{aligned}
\int_{t}^{T} L(-\dot{\eta}, \eta) \mathrm{d} \tau & =\frac{1}{4} \int_{t}^{T} a^{i j}(\eta(\tau))\left(-\dot{\eta}_{i}(\tau)+b_{i}(\eta(\tau))\right)\left(-\dot{\eta}_{j}(\tau)+b_{j}(\eta(\tau))\right) \mathrm{d} \tau \\
& \leq 4 \int_{t}^{T}\left(\bar{\Gamma}(\tau) \zeta^{\prime}(\tau)\right) \cdot\left(\bar{Q}(\tau) \bar{\Gamma}(\tau) \zeta^{\prime}(\tau)\right) \mathrm{d} \tau+C T\left|z^{\prime}\right|^{3}
\end{aligned}
$$

In view of (7.7) and (7.2), we have

$$
\begin{aligned}
\int_{t}^{T} L(-\dot{\eta}, \eta) \mathrm{d} \tau & \leq \int_{t}^{T}\left(\bar{\Gamma} \zeta^{\prime}\right) \cdot\left(\bar{B} \zeta^{\prime}-\dot{\zeta}^{\prime}\right) \mathrm{d} \tau+C T\left|z^{\prime}\right|^{3} \\
& =-\int_{t}^{T} \frac{\mathrm{~d}}{\mathrm{~d} \tau}\left(\bar{\Gamma} \zeta^{\prime} \cdot \zeta^{\prime}\right) \mathrm{d} \tau+\int_{t}^{T}\left(\bar{\Gamma} \zeta^{\prime} \cdot\left(\bar{B} \zeta^{\prime}+\dot{\zeta}^{\prime}\right)+\dot{\bar{\Gamma}} \zeta^{\prime} \cdot \zeta^{\prime}\right) \mathrm{d} \tau+C T\left|z^{\prime}\right|^{3} \\
& \left.\leq \bar{\Gamma}(t) z^{\prime} \cdot z^{\prime}+\int_{t}^{T}\left(\dot{\bar{\Gamma}}-4 \bar{\Gamma} \bar{Q} \bar{\Gamma}+\bar{\Gamma} \bar{B}+\bar{B}^{*} \bar{\Gamma}\right) \zeta^{\prime} \cdot \zeta^{\prime}\right) \mathrm{d} \tau+C T\left|z^{\prime}\right|^{3} \\
& =\bar{\Gamma}(t) z^{\prime} \cdot z^{\prime}+C T\left|z^{\prime}\right|^{3}
\end{aligned}
$$

If we choose $T:=\frac{1}{\delta} \log \frac{1}{\left|z^{\prime}\right|}$, then $\operatorname{dist}(\eta(T), \mathcal{C})=O\left(\left|z^{\prime}\right|^{2}\right)$ and

$$
\begin{equation*}
\int_{t}^{T} L(-\dot{\eta}, \eta) \mathrm{d} \tau \leq \bar{\Gamma}(t) z^{\prime} \cdot z^{\prime}+C\left|z^{\prime}\right|^{3} \log \frac{1}{\left|z^{\prime}\right|} \tag{7.9}
\end{equation*}
$$

For constructing $\zeta_{N}$, we will need the following facts. From the definition of $\mathcal{T}_{i j}$ in (7.5), it follows that $\mathcal{T}_{i N}(\tau)=\dot{\xi}_{i}(\tau) /|\dot{\xi}(\tau)|$. Then, since $\ddot{\xi}_{i}(\tau)=\frac{\partial b_{i}}{\partial x_{j}}(\xi(\tau)) \dot{\xi}_{j}(\tau)$, we have

$$
\frac{\partial b_{i}}{\partial x_{j}}(\xi(\tau)) \mathcal{T}_{j N}(\tau)-\dot{\mathcal{T}}_{i N}(\tau)=\frac{\dot{\xi}_{i}(\tau) \dot{\xi}_{j}(\tau) \ddot{\xi}_{j}(\tau)}{|\dot{\xi}(\tau)|^{3}}=\mathcal{T}_{i N}(\tau) \frac{\dot{\xi}_{j}(\tau) \ddot{\xi}_{j}(\tau)}{|\dot{\xi}(\tau)|^{2}}, \quad \forall i=1, \ldots, N
$$

Multiplying this by $\mathcal{T}^{-1}(\tau)$, we conclude that

$$
\mathcal{T}_{k i}^{-1}(\tau)\left(\frac{\partial b_{i}}{\partial x_{j}}(\xi(\tau)) \mathcal{T}_{j N}(\tau)-\dot{\mathcal{T}}_{i N}(\tau)\right)=\left\{\begin{array}{l}
\dot{\xi}_{j}(\tau) \ddot{\xi}_{j}(\tau) /|\dot{\xi}(\tau)|^{2}, \quad \text { if } k=N,  \tag{7.10}\\
0 \quad \text { if } k<N
\end{array}\right.
$$

We proceed with constructing $\zeta_{N}$. By the definition of $\eta$, we have

$$
\begin{aligned}
\dot{\eta}_{i}(\tau)- & b_{i}(\eta(\tau))=\frac{\mathrm{d}}{\mathrm{~d} t}\left(\xi_{i}(\tau)+\mathcal{T}_{i k}(\tau) \zeta_{k}^{\prime}(\tau)+\mathcal{T}_{i N}(\tau) \zeta_{N}(\tau)\right)-b_{i}(\eta(\tau))+O\left(\left|z^{\prime}\right|^{2}\right) \\
= & \left(\dot{\mathcal{T}}_{i k}(\tau)-\frac{\partial b_{i}}{\partial x_{j}}(\xi(\tau)) \mathcal{T}_{j k}(\tau)\right) \zeta_{k}^{\prime}(\tau)+\left(\dot{\mathcal{T}}_{i N}(\tau)-\frac{\partial b_{i}}{\partial x_{j}}(\xi(\tau)) \mathcal{T}_{j N}(\tau)\right) \zeta_{N}(\tau) \\
& +\mathcal{T}_{i k}(\tau) \dot{\zeta}_{k}^{\prime}(\tau)+\mathcal{T}_{i N}(\tau) \dot{\zeta}_{N}(\tau)+O\left(\left|z^{\prime}\right|^{2}\right) \\
& =\mathcal{T}_{i k}(\tau) \mathcal{T}_{k r}^{-1}(\tau)\left\{\left(\dot{\mathcal{T}}_{r k}(\tau)-\frac{\partial b_{r}}{\partial x_{j}}(\xi(\tau)) \mathcal{T}_{j k}(\tau)\right) \zeta_{k}^{\prime}(\tau)\right.
\end{aligned}
$$

$$
\begin{align*}
& \left.+\left(\dot{\mathcal{T}}_{r N}(\tau)-\frac{\partial b_{r}}{\partial x_{j}}(\xi(\tau)) \mathcal{T}_{j N}(\tau)\right) \zeta_{N}(\tau)\right\} \\
& +\mathcal{T}_{i k}(\tau) \dot{\zeta}_{k}^{\prime}(\tau)+\mathcal{T}_{i N}(\tau) \dot{\zeta}_{N}(\tau)+O\left(\left|z^{\prime}\right|^{2}\right) \tag{7.11}
\end{align*}
$$

Substituting for $\dot{\zeta}^{\prime}$ the expression on the right-hand side of (7.7) and considering (7.10) yields
$b_{i}(\eta(\tau))-\dot{\eta}_{i}(\tau)=4 a_{i j}(\xi(\tau)) \mathcal{T}_{l j}^{-1}(0, \xi(\tau)) \bar{\Gamma}_{l m}(\tau) \zeta_{m}^{\prime}(\tau)+O\left(\left|z^{\prime}\right|^{2}\right)-\mathcal{T}_{i N}(\tau) \dot{\zeta}_{N}(\tau)+\mathcal{T}_{i N}(\tau) R(\tau)$,
where

$$
\begin{aligned}
R(\tau)= & \mathcal{T}_{N i}^{-1}(\tau) \frac{\partial b_{i}}{\partial x_{j}}(\xi(\tau))\left(\mathcal{T}_{j N}(\tau) \zeta_{N}(\tau)+\mathcal{T}_{j l}(\tau) \zeta_{l}^{\prime}(\tau)\right) \\
& -\mathcal{T}_{N i}^{-1}(\tau)\left(\dot{\mathcal{T}}_{i N}(\tau) \zeta_{N}(\tau)+\dot{\mathcal{T}}_{i l}(\tau) \zeta_{l}^{\prime}(\tau)\right)-4 \mathcal{T}_{N i}^{-1}(\tau) a_{i j}(\xi(\tau)) \mathcal{T}_{l j}^{-1}(\tau) \bar{\Gamma}_{l m}(\tau) \zeta_{m}^{\prime}(\tau)
\end{aligned}
$$

Thus, in order to make (7.8) hold, we choose $\zeta_{N}(\tau)$ as a solution of the following equation

$$
\begin{equation*}
\dot{\zeta}_{N}(\tau)=R(\tau) \tag{7.12}
\end{equation*}
$$

with the initial condition $\zeta_{N}(t)=0$.
From (7.10) it follows that $|\dot{\xi}(\tau)|$ solves

$$
\frac{\mathrm{d}}{\mathrm{~d} t}|\dot{\xi}(\tau)|=\left(\mathcal{T}_{N i}^{-1}(\tau) \frac{\partial b_{i}}{\partial x_{j}}(\xi(\tau)) \mathcal{T}_{j N}(\tau)-\mathcal{T}_{N i}^{-1}(\tau) \dot{\mathcal{T}}_{i N}(\tau)\right)|\dot{\xi}(\tau)|
$$

and we can write the solution $\zeta_{N}(\tau)$ of (7.12) as

$$
\begin{align*}
\zeta_{N}(\tau) & =|\dot{\xi}(\tau)| \int_{0}^{\tau}\left(\mathcal{T}_{N i}^{-1}(s) \frac{\partial b_{i}}{\partial x_{j}}(\xi(s)) \mathcal{T}_{j l}(s) \zeta_{l}^{\prime}(s)-\mathcal{T}_{N i}^{-1}(s) \dot{\mathcal{T}}_{i l}(s) \zeta_{l}^{\prime}(s)\right) \frac{\mathrm{d} s}{|\dot{\xi}(s)|} \\
& -4|\dot{\xi}(\tau)| \int_{0}^{\tau} \mathcal{T}_{N i}^{-1}(s) a_{i j}(\xi(s)) \mathcal{T}_{l j}^{-1}(s) \bar{\Gamma}_{l m}(s) \zeta_{m}^{\prime}(s) \frac{\mathrm{d} s}{|\dot{\xi}(s)|} \tag{7.13}
\end{align*}
$$

From (7.13) we derive the uniform bound $\left|\zeta_{N}(\tau)\right| \leq C\left|z^{\prime}\right|$.
It remains to construct $\eta(\cdot)$ for $\tau>T$ in such a way that it reaches the cycle in a finite time. To this end, we set

$$
\begin{equation*}
\eta(\tau)=X\left(\zeta^{\prime}(T)(T+1-\tau), z_{N}(\xi(\tau))+\zeta_{N}(T)|\dot{\xi}(\tau)| / / \dot{\xi}(T) \mid\right), \quad T \leq \tau \leq T+1 \tag{7.14}
\end{equation*}
$$

Then for every $i=1, \ldots, N$,

$$
\begin{aligned}
\dot{\eta}_{i}(\tau) & =\frac{\mathrm{d}}{\mathrm{~d} \tau}\left(\xi_{i}(\tau)+\mathcal{T}_{i N}(\tau) \zeta_{N}(T)|\dot{\xi}(\tau)| /|\dot{\xi}(T)|\right)+O\left(\left|z^{\prime}\right|^{2}\right) \\
& =\dot{\xi}_{i}(\tau)+\ddot{\xi}_{i}(\tau) \zeta_{N}(T) /|\dot{\xi}(T)|+O\left(\left|z^{\prime}\right|^{2}\right) \\
& =b_{i}(\xi(\tau))+\frac{\partial b_{i}}{\partial x_{j}}(\xi(\tau)) \dot{\xi}_{j}(\tau) \zeta_{N}(T) /|\dot{\xi}(T)|+O\left(\left|z^{\prime}\right|^{2}\right)=b_{i}(\eta(\tau))+O\left(\left|z^{\prime}\right|^{2}\right)
\end{aligned}
$$

Therefore, the following bound holds

$$
\begin{equation*}
\int_{T}^{T+1} a^{i j}(\eta(\tau))\left(-\dot{\eta}_{i}+b_{i}(\eta)\right)\left(-\dot{\eta}_{j}+b_{j}(\eta)\right) \mathrm{d} \tau \leq C\left|z^{\prime}\right|^{4} . \tag{7.15}
\end{equation*}
$$

Combining the last relation with (7.9) yields (7.6).
In order to construct sub- and super-solution of (6.1), we consider the solution $\bar{D}(t)$ of the matrix equation

$$
\begin{equation*}
\dot{\bar{D}}+\bar{D}(\bar{B}-4 \bar{Q} \bar{\Gamma})+(\bar{B}-4 \bar{Q} \bar{\Gamma})^{*} \bar{D}=-2 I, \tag{7.16}
\end{equation*}
$$

given by

$$
\begin{equation*}
\bar{D}(t)=2 \int_{t}^{\infty} \bar{\Psi}^{*}(\tau, t) \bar{\Psi}(\tau, t) \mathrm{d} \tau \tag{7.17}
\end{equation*}
$$

where $\bar{\Psi}(\tau, t)$ is the fundamental matrix solution of

$$
\frac{\partial \bar{\Psi}}{\partial \tau}=(\bar{B}(\tau)-4 \bar{Q}(\tau) \bar{\Gamma}(\tau)) \bar{\Psi}, \quad \bar{\Psi}(t, t)=I .
$$

As already mentioned in the proof of Lemma 9 , this solution $\bar{\Psi}(\tau, t)$ decays exponentially as $\tau \rightarrow+\infty$ and therefore the integral in (7.17) converges. Then it defines a $P$-periodic positive symmetric solution of (7.16). It follows from (7.2) and (7.16) that $\bar{\Gamma}_{\delta}^{ \pm}:=\bar{\Gamma} \pm \delta \bar{D}$ satisfy for sufficiently small $\delta>0$
$4 \bar{\Gamma}_{\delta}^{+} \bar{Q} \bar{\Gamma}_{\delta}^{+}-\frac{\mathrm{d}}{\mathrm{d} t} \bar{\Gamma}_{\delta}^{+}-\bar{\Gamma}_{\delta}^{+} \bar{B}-\bar{B}^{*} \bar{\Gamma}_{\delta}^{+} \geq \delta I$ and $4 \bar{\Gamma}_{\delta}^{-} \bar{Q} \bar{\Gamma}_{\delta}^{-}-\frac{\mathrm{d}}{\mathrm{d} t} \bar{\Gamma}_{\delta}^{-}-\bar{\Gamma}_{\delta}^{+} \bar{B}-\bar{B}^{*} \bar{\Gamma}_{\delta}^{-} \leq-\delta I$.
Now define the functions $W_{\delta}^{ \pm}(z)$ by

$$
W_{\delta}^{ \pm}(z)=\left(\bar{\Gamma}_{\delta}^{ \pm}(t)\right)_{i j} z_{i}^{\prime} z_{j}^{\prime} \quad \text { where } z_{N}=z_{N}(\xi(t))
$$

these functions satisfy
$S\left(\nabla_{z} W_{\delta}^{-}(z), z\right) \leq-\frac{\delta}{2}\left|z^{\prime}\right|^{2} \quad$ and $\quad S\left(\nabla_{z} W_{\delta}^{+}(z), z\right) \geq \frac{\delta}{2}\left|z^{\prime}\right|^{2} \quad$ for sufficiently small $\left|z^{\prime}\right|$.
The latter inequalities follow directly form the definitions of $W_{\delta}^{-}(z)$ and $W_{\delta}^{+}(z)$.
Lemma 10 For sufficiently small $\delta>0$, the strict pointwise inequalities $W_{\delta}^{-}(z)<$ $W(X(z))<W_{\delta}^{+}(z)$ hold for $z^{\prime}$ from a punctured neighbourhood of zero.

Proof The first inequality $W_{\delta}^{-}(z)<W(X(z))$ can be proved similarly to Lemma 19 in [12] (see also the proof of Lemma 8), using (7.18). The second inequality $W(X(z))<W_{\delta}^{+}(z)$ follows immediately from Lemma 9.

At this point, we have constructed functions $W_{\delta}^{ \pm}$satisfying conditions of Lemma 5. Next, we define the test functions $W_{\delta, \varepsilon}^{ \pm}$by

$$
W_{\delta, \varepsilon}^{ \pm}:=W_{\delta}^{ \pm}-\varepsilon \bar{\Phi}_{\delta}^{ \pm}(t)
$$

where $z_{N}$ and $t$ are related by $z_{N}=z_{N}(\xi(t))$, and $\bar{\Phi}_{\delta}^{ \pm}(t)$ are periodic solutions of the ODEs

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} t} \bar{\Phi}_{\delta}^{ \pm}(t)=-2 \operatorname{tr}\left(\bar{Q}(t) \bar{\Gamma}_{\delta}^{ \pm}(t)\right)+c(\xi(t))+\frac{2}{P} \int_{0}^{P} \operatorname{tr}\left(\bar{Q}(\tau) \bar{\Gamma}_{\delta}^{ \pm}(\tau)\right) \mathrm{d} \tau-\frac{1}{P} \int_{0}^{P} c(\xi(\tau)) \mathrm{d} \tau \tag{7.19}
\end{equation*}
$$

The first two terms on the right-hand side here are introduced in order to compensate the discrepancy of order $\varepsilon$ in Equation (6.2). Indeed, the test functions $W_{\delta, \varepsilon}^{ \pm}$constructed in this way satisfy for sufficiently small $\left|z^{\prime}\right|$

$$
\begin{aligned}
& \pm a_{i j}(X(z)) \alpha_{k i}(z) \frac{\partial}{\partial z_{k}}\left(\alpha_{l j}(z) \frac{\partial W_{\delta, \varepsilon}^{ \pm}}{\partial z_{l}}\right) \mp \frac{1}{\varepsilon} S\left(\nabla_{z} W_{\delta, \varepsilon}^{ \pm}, z\right) \mp c(X(z)) \\
& \leq \frac{1}{P} \int_{0}^{P} c(\xi(\tau)) \mathrm{d} \tau-\frac{2}{P} \int_{0}^{P} \operatorname{tr}\left(\bar{Q}(\tau) \bar{\Gamma}_{\delta}^{ \pm}(\tau)\right) \mathrm{d} \tau+O\left(\varepsilon+\left|z^{\prime}\right|\right)
\end{aligned}
$$

In order to complete the proof of the fact that $W_{\delta, \varepsilon}^{ \pm}$satisfy (4.4) and (4.5), it remains to observe that $\int_{0}^{P} \operatorname{tr}\left(\bar{Q}(\tau) \bar{\Gamma}_{\delta}^{ \pm}(\tau)\right) \mathrm{d} \tau \rightarrow \int_{0}^{P} \operatorname{tr}(\bar{Q}(\tau) \bar{\Gamma}(\tau)) \mathrm{d} \tau$ as $\delta \rightarrow 0$ and use the identity

$$
2 \int_{0}^{P} \operatorname{tr}(\bar{Q}(\tau) \bar{\Gamma}(\tau)) \mathrm{d} \tau=\sum_{\Theta_{i}>1} \log \Theta_{i}
$$

(see Proposition 5.1 in [13]), where $\Theta_{i}$ are absolute values of eigenvalues of the linearized Poincaré map (corresponding to the ODE $\dot{x}=b(x)$ near $\mathcal{C}$ ).

## 8. Construction of test functions: case of limit cycles on $\partial \Omega$

In the case when $\mathrm{ODE} \dot{x}=b_{\tau}(x)$ on $\partial \Omega$ has a limit cycle $\mathcal{C}$ which is significant component of the Aubry set, the analysis combines the ideas of Sections 6 and 7. We pass to the local coordinates in a neighbourhood of $\mathcal{C}$ via a map $x=X\left(z_{1}, \ldots, z_{N-1}, z_{N}\right)$, where $z_{N}=z_{N}(x)$ is the distance from $x$ to $\partial \Omega$ (positive for $x \in \Omega$ ) and $\left(z_{1}, \ldots, z_{N-1}\right)$ are coordinates on $\partial \Omega$. The coordinate $z_{N-1}(x)$ represents the arc length parametrization on $\mathcal{C}$ and other coordinates $z^{\prime}=\left(z_{1}, \ldots, z_{N-2}\right)$ are chosen so that the map $X\left(z^{\prime}, z_{N-1}, z_{N}\right)$ is $C^{2}$-smooth; moreover, $z^{\prime}=0$ when $x \in \mathcal{C}$, and $\left(\frac{\partial X_{i}}{\partial z_{j}}(z)\right)_{i, j=\overline{1, N}}$ is an orthogonal matrix when $z_{N}=0$ and $z^{\prime}=0$ (on the cycle). This change of coordinates leads to equations of the form (6.1) and (6.2) for $W(X(z))$ and $W_{\varepsilon}(X(z))$.

We use the following ansatz for $W$,

$$
\begin{equation*}
W(X(z))=\widehat{\Gamma}_{i j}(t) z_{i}^{\prime} z_{j}^{\prime}+o\left(\left|z^{\prime}\right|^{2}\right) \tag{8.1}
\end{equation*}
$$

where $\widehat{\Gamma}$ is now $(N-2) \times(N-2)$ symmetric $P$-periodic matrix ( $P$ being the period of the cycle $\mathcal{C})$, and $t$ refers to the parametrization $t \rightarrow \xi(t)$ of $\mathcal{C}$ such that $\dot{\xi}(t)=b_{\tau}(\xi(t))$. Moreover, $\widehat{\Gamma}$ is chosen to be the maximal $P$-periodic solution of the Riccati matrix equation

$$
\frac{\mathrm{d}}{\mathrm{~d} t} \widehat{\Gamma}=4 \widehat{\Gamma} \widehat{Q} \widehat{\Gamma}-\widehat{\Gamma} \widehat{B}-\widehat{B}^{*} \widehat{\Gamma}
$$

with $(N-2) \times(N-2)$ matrices $\widehat{Q}(t)$ and $\widehat{B}(t)$ whose entries are given by the same formulas as (7.3) and (7.4).

Lemma 11 For sufficiently small $\left|z^{\prime}\right|$ and $\left|z_{N}\right|$ the following bound holds uniformly in $t \in[0, P)$

$$
\begin{equation*}
W\left(X\left(z^{\prime}, z_{N-1}(\xi(t)), z_{N}\right) \leq \widehat{\Gamma}_{i j}(t) z_{i}^{\prime} z_{j}^{\prime}+C\left(\left|z^{\prime}\right|^{2} \log \frac{1}{\left|z^{\prime}\right|}+\left|z_{N}\right|\left|z^{\prime}\right|^{2}+\left|z_{N}\right|^{3}\right)\right. \tag{8.2}
\end{equation*}
$$

Proof First consider the case $z_{N}=0$. As in Lemma 9, we use representation (2.9) and consider the solution $\zeta^{\prime}(\tau)$ of the $\operatorname{ODE} \dot{\zeta}^{\prime}(\tau)=(\widehat{B}-4 \widehat{Q} \widehat{\Gamma}) \zeta(\tau)$ for $\tau>t$ with the initial condition $\zeta^{\prime}(t)=z^{\prime}$. It decays exponentially as $\tau \rightarrow \infty,\left|\zeta^{\prime}\right| \leq C e^{-\delta \tau}\left|z^{\prime}\right|$ for some $\delta>0$. Next we introduce $\zeta_{N-1}$ analogously to $\zeta_{N}$ introduced in Lemma 9, i.e. $\zeta_{N-1}$ solves

$$
\begin{aligned}
\dot{\zeta}_{N-1}(\tau)= & \mathcal{T}_{(N-1) i}^{-1}\left(0, z_{N-1}(\xi(\tau), 0) \frac{\partial b_{i}}{\partial x_{j}}(\xi(\tau))\left(\mathcal{T}_{j(N-1)}(\tau) \zeta_{N-1}(\tau)+\mathcal{T}_{j l}(\tau) \zeta_{l}^{\prime}(\tau)\right)\right. \\
& -\mathcal{T}_{(N-1) i}^{-1}\left(0, z_{N-1}(\xi(\tau)), 0\right)\left(\dot{\mathcal{T}}_{i(N-1)}(\tau) \zeta_{N-1}(\tau)+\dot{\mathcal{T}}_{l l}(\tau) \zeta_{l}^{\prime}(\tau)\right) \\
& -4 \mathcal{T}_{(N-1) i}^{-1}\left(0, z_{N-1}(\xi(\tau)), 0\right) a_{i j}(\xi(\tau)) \mathcal{T}_{l j}^{-1}\left(0, z_{N-1}(\xi(\tau)), 0\right) \widehat{\Gamma}_{l m}(\tau) \zeta_{m}^{\prime}(\tau), \quad \tau>t
\end{aligned}
$$

where $\left(\mathcal{T}_{i j}^{-1}(z)\right)_{i, j=\overline{1, N}}$ is the matrix inverse to $\left(\frac{\partial X_{i}}{\partial z_{j}}(z)\right)_{i, j=\overline{1, N}}$ and $\mathcal{T}_{i j}(\tau)=$ $\frac{\partial X_{i}}{\partial z_{j}}\left(0, z_{N-1}(\xi(\tau)), 0\right)$.

Finally, we define $\eta(\tau)$ by
$\eta(\tau)=\left\{\begin{array}{l}X\left(\zeta^{\prime}(\tau), z_{N-1}(\xi(\tau))+\zeta_{N-1}(\tau), 0\right), \quad t \leq \tau<T \\ X\left(\zeta^{\prime}(T)(T+1-\tau), z_{N-1}(\xi(\tau))+\zeta_{N-1}(T)|\dot{\xi}(\tau)| /|\dot{\xi}(T)|, 0\right), \quad T \leq \tau<T+1,\end{array}\right.$
with $T:=\frac{1}{\delta} \log \frac{1}{\left|z^{\prime}\right|}$, and the control $v(\tau)$ by

$$
v(\tau)=\left\{\begin{array}{l}
\dot{\eta}+v(\eta)\left(b_{v}(\eta)+R_{1}(\tau)\right), \quad t \leq \tau<T \\
\dot{\eta}+v(\eta) b_{v}(\eta), \quad T \leq \tau \leq T+1
\end{array}\right.
$$

where $z_{N-1}=z_{N-1}(\xi(\tau))$, and

$$
\begin{aligned}
R_{1}(\tau)= & -\mathcal{T}_{(N-1) i}^{-1}\left(0, z_{N-1}, 0\right)\left(\dot{\mathcal{T}}_{i(N-1)}(\tau) \zeta_{N-1}(\tau)+\dot{\mathcal{T}}_{i l}(\tau) \zeta_{l}^{\prime}(\tau)\right. \\
& \left.+4 a_{i j}(\xi(\tau)) \mathcal{T}_{l j}^{-1}\left(0, z_{N-1}, 0\right) \widehat{\Gamma}_{l m} \zeta_{m}^{\prime}\right)
\end{aligned}
$$

Letting

$$
\alpha(\tau):=\left\{\begin{array}{l}
b_{\nu}(\eta)+R_{1}(\tau), \quad t \leq \tau<T \\
b_{v}(\eta), \quad T \leq \tau \leq T+1,
\end{array}\right.
$$

observe that for this control $v(\tau)$, the pair $(\eta(\tau), \alpha(\tau))$ solves (2.8) on $(t, T+1)$ with the initial value $\eta(t)=x\left(=X\left(z^{\prime}, \xi(t), 0\right)\right)$, as far as $\alpha(\tau) \geq 0$ for all $\tau \in(t, T+1)$. Since $b_{\nu}(\xi(\tau))>0$, the latter condition is satisfied, provided that $\left|z^{\prime}\right|$ is sufficiently small. Then the proof of (8.2) follows exactly the line of Lemma 9.

In the case when $z_{N}(x)>0$, we construct a curve $\eta(\tau)$ connecting $x$ with a point $y$ on $\partial \Omega$ by setting
$\eta(\tau)=X\left(z^{\prime}, z_{N-1}(\xi(t)), z_{N}+b_{\nu}(\xi(t))(t-\tau)\right)$ for all $\tau \geq t$ such that $z_{N}+b_{\nu}(\xi(t))(t-\tau) \geq 0$.

Let $t+\Delta t$ be the time when $\eta(\tau)$ reaches $\partial \Omega$ (at the point $y=\eta(t+\Delta t))$ then $\Delta t=O\left(z_{N}\right)$. It follows from the construction of $\eta(\tau)$ that

$$
\int_{t}^{t+\Delta t} a^{i j}(\eta(\tau))\left(-\dot{\eta}_{i}+b_{i}(\tau)\right)\left(-\dot{\eta}_{j}+b_{j}(\eta)\right) \mathrm{d} \tau \leq C\left(\left|z^{\prime}\right|^{2}+z_{N}^{2}\right) z_{N}
$$

Then extending $\eta(\tau)$ along $\partial \Omega$ as described above we complete the proof of the Lemma.

Now, we construct test functions $W_{\delta}^{ \pm}\left(z^{\prime}, z_{N-1}(\xi(t)), z_{N}\right):=(\widehat{\Gamma} \pm \delta \widehat{D})_{i j}(t) z_{i}^{\prime} z_{j}^{\prime} \pm \delta z_{N}^{2}$ for (sufficiently small) $\delta>0$, where the $P$-periodic symmetric matrix $\hat{D}(t)>0$ is defined analogously to (7.17). Then

$$
S\left(\nabla_{z} W_{\delta}^{-}, z\right) \leq-\delta\left(\left|z^{\prime}\right|^{2}+b_{v}(\xi(t)) z_{N}\right)
$$

for sufficiently small $\left|z^{\prime}\right|$ and $z_{N} \geq 0$. This yields the following bound

$$
W_{\delta}^{-}<W(X(z)) \text { for sufficiently small }\left|z^{\prime}\right| \text { and } z_{N}\left(\text { when }\left|z^{\prime}\right|+z_{N}>0\right)
$$

whose proof is analogous to that of the lower bound in Lemma 8. Thus functions $W_{\delta}^{ \pm}$satisfy the conditions of Lemma 5. Finally, we define the test functions $W_{\delta, \varepsilon}^{ \pm}$by

$$
W_{\delta, \varepsilon}^{ \pm}:=W_{\delta}^{ \pm} \pm \varepsilon \widehat{\Phi}_{\delta}^{ \pm} \mp \varepsilon^{2} z_{N}, \text { where } z_{N-1}=z_{N-1}(\xi(t)),
$$

with $\widehat{\Phi}_{\delta}^{ \pm}$being solutions of

$$
\frac{\mathrm{d}}{\mathrm{~d} t} \widehat{\Phi}_{\delta}^{ \pm}(t)=-2 \operatorname{tr}\left(\widehat{Q}(t) \widehat{\Gamma}_{\delta}^{ \pm}(t)\right)+c(\xi(t))+\frac{2}{P} \int_{0}^{P} \operatorname{tr}\left(\widehat{Q}(\tau) \widehat{\Gamma}_{\delta}^{ \pm}(\tau)\right) \mathrm{d} \tau-\frac{1}{P} \int_{0}^{P} c(\xi(\tau)) \mathrm{d} \tau
$$

These functions $W_{\delta, \varepsilon}^{ \pm}$satisfy the conditions of Lemma 5 .

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