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# Spectrum of the Second Variation 

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#### Abstract

Second variation of a smooth optimal control problem at a regular extremal is a symmetric Fredholm operator. We study asymptotics of the spectrum of this operator and give an explicit expression for its determinant in terms of solutions of the Jacobi equation. In the case of the least action principle for the harmonic oscillator we obtain a classical Euler identity $\prod_{n=1}^{\infty}\left(1-\frac{x^{2}}{(\pi n)^{2}}\right)=\frac{\sin x}{x}$. General case may serve as a rich source of new nice identities.


## 1 Introduction

This research was initially motivated by the analysis of asymptotic expansions related to semigroups generated by hypo-elliptic operators. These expansions can be often interpreted as infinite-dimensional versions of the standard asymptotics for expressions of the form $\int_{\mathbb{R}^{N}} a(u) e^{\frac{\varphi(u)}{-t}} d u$ or $\int_{\mathbb{R}^{N}} a(u) e^{\frac{\varphi(u)}{i t}} d u$ as $t \rightarrow 0$. It is well-known that the asymptotics is localized in the critical points $x_{i}$ of $\varphi$, and the quantities $\left(\operatorname{det} D_{x_{i}}^{2}\right)^{-\frac{1}{2}}$ play a crucial role in the expansions. In the infinite-dimensional version we deal with a "path integral", and $\varphi$ is an "action functional" whose critical points are extremals of the appropriate optimal control problem. So we try to analyse and compute $\operatorname{det} D_{x_{i}}^{2}$ in this case. What we obtain, can be considered as a modification and generalization of the classical "Hill's determinant" in Rational Mechanics (see [2]).

Before the calculation of the determinant, we study asymptotics of the spectrum of the second variation $D_{u}^{2} \varphi$ that is a symmetric Fredholm operator of the form $I+K$, where $K$ is a compact Hilbert-Schmidt operator. The

[^0]point is that $K$ is usually NOT a trace class operator so that the trace of $K$ and the determinant of $I+K$ are not well-defined in the standard sense.

Anyway, a specific structure of the spectral asymptotics allows to define and compute everything. This asymptotic is described in Theorem 1. Actually, there is an important symmetry in the asymptotics that leads to a cancellation of slow convergent to zero positive and negative terms when we compute the sum or the product.

There is a general reason for this symmetry. It concerns an evolutionary nature of the optimal control problems: the space of available control functions grows with time. This evolutionary structure results in the following property (I would call it "causality") of the operator $K$ : there exists a Volterra operator $V$ on $L^{2}\left([0,1] ; \mathbb{R}^{m}\right)$ and a finite codimension subspace of $L^{2}\left([0,1] ; \mathbb{R}^{m}\right)$ such that $\langle K u, u\rangle=\langle V u, u\rangle$ for any $u$ from the subspace. In other words, our symmetric operator imitates a triangle one on a subspace of finite codimension. We expect that all causal operators have spectral asymptotics similar to the asymptotics from Theorem 1, although we do not study general causal operators in this paper.

In Theorem 2 we give an explicit expression for $\operatorname{det}(I+K)$ in terms of solutions of the Jacobi equation that generalize classical Jacobi fields. We also give an explicit integral expression for $\operatorname{tr} K$; moreover, integral formulas for all elementary symmetric functions of the eigenvalues of $K$ can be recovered from the main determinantal formula.

I hope that this result would serve as an effective summation method, a way to get explicit expressions for infinite products and sums of interesting series. A simple example: for the 1-dimensional linear control system $\dot{x}=a x+u$ with the quadratic cost $\varphi(u)=\int_{0}^{1} u^{2}(t)-\left(a^{2}+b^{2}\right) x^{2}(t) d t$ our determinantal identity reads: $\prod_{n=1}^{\infty}\left(1-\frac{a^{2}+b^{2}}{a^{2}+(\pi n)^{2}}\right)=\frac{a \sin b}{b \operatorname{sh} a}$; the case $a=0$ corresponds to the famous Euler identity $\prod_{n=1}^{\infty}\left(1-\frac{b^{2}}{(\pi n)^{2}}\right)=\frac{\sin b}{b}$.

## 2 Preliminaries

We consider a smooth control system of the form:

$$
\begin{equation*}
\dot{q}=f_{u}^{t}(q), \quad q \in M, u \in U \tag{1}
\end{equation*}
$$

with a fixed initial point $q_{0} \in M$. Here $M$ and $U$ are smooth manifolds (without border) and the vector $f_{u}^{t}(q) \in T_{q} M$ smoothly depends on $(q, u) \in$ $M \times U$ and is measurable bounded with respect to $t \in[0,1]$.

Let $u(\cdot):[0,1] \rightarrow U$ be a measurable map such that $u([0,1])$ is contained in a compact subset of $U$. The Cauchy problem

$$
\dot{q}=f_{u(t)}^{t}(q), \quad q(0)=q_{0},
$$

has a unique Lipschitzian solution $t \mapsto q(t ; u(\cdot))$ defined on an interval in $\mathbb{R}$. We say that $u(\cdot)$ is an admissible control and $q(\cdot ; u(\cdot))$ is a corresponding admissible trajectory if the domain of this solution contains $[0,1]$. We denote the set of all admissible controls by $\mathcal{U}$; then $\mathcal{U}$ is an open subset of $L^{\infty}([0,1] ; U)$. Hence $\mathcal{U}$ is a smooth Banach manifold modelled on the space $L^{\infty}\left([0,1] ; \mathbb{R}^{\operatorname{dim} U}\right)$.

Given $t \in[0,1]$ we define the "evaluation map" $F_{t}: \mathcal{U} \rightarrow M$ by the formula $F_{t}(u(\cdot))=q(t ; u(\cdot))$; then $F_{t}$ is a smooth map from the Banach manifold $\mathcal{U}$ to $M$.

Let $\ell: M \times U \rightarrow \mathbb{R}$ be a smooth "Lagrangian". We consider functionals $\varphi_{t}: \mathcal{U} \rightarrow \mathbb{R}, 0 \leq t \leq 1$, defined by the formula:

$$
\begin{equation*}
\varphi_{t}(u(\cdot))=\int_{0}^{t} \ell(q(\tau ; u(\cdot)), u(\tau)) d \tau \tag{2}
\end{equation*}
$$

Definition 1. We say that $u \in \mathcal{U}$ is a normal extremal control if there exists $\lambda_{1} \in T_{F_{1}(u)}^{*} M$ such that $\lambda_{1} D_{u} F_{1}=d_{u} \varphi_{1}$; here $\lambda_{1} D_{u} F_{1}$ is the composition of $D_{u} F_{1}: T_{u} \mathcal{U} \rightarrow T_{F_{1}(u)} M$ and $\lambda_{1}: T_{F_{1}(u)} M \rightarrow \mathbb{R}$. We say that a normal extremal control is strictly normal if it is a regular point of $F_{1}$.

Strictly normal extremal controls are just critical points of $\left.\varphi\right|_{F_{1}^{-1}(q)}, q \in$ $M$. Let $t \in[0,1]$; it is easy to see that the restriction $\left.u\right|_{[0, t]}$ of a normal extremal control $u$ is also a normal extremal control in the following sense:

$$
\begin{equation*}
\exists \lambda_{t} \in T_{F_{t}(u)^{*} M}^{M} \text { such that } \lambda_{t} D_{u} F_{t}=d_{u} \varphi_{t} . \tag{3}
\end{equation*}
$$

The "differentiation of the identity (3) with respect to $t$ " leads to the following Hamiltonian characterization of normal extremal controls.

A family of Hamiltonians $h_{u}^{t}: T^{*} M \rightarrow \mathbb{R}, u \in U$, is defined by the formula:

$$
h_{u}^{t}(\lambda)=\left\langle\lambda, f_{u}^{t}(q)\right\rangle-\ell(q, u), \quad g \in M, \lambda \in T_{q}^{*} M .
$$

Let $\sigma$ be the canonical symplectic form on $T^{*} M$ and $\pi: T^{*} M \rightarrow M$ be the standard projection, $\pi\left(T_{q}^{*} M\right)=q$. Recall that $\sigma=d s$ where $s$ is a Liouville (or tautological) 1-form, $\left\langle s_{\lambda}, \eta\right\rangle=\left\langle\lambda, \pi_{*} \eta\right\rangle, \forall \lambda \in T^{*} M, \eta \in T_{\lambda}\left(T^{*} M\right)$. Given a smooth function $h: T^{*} M \rightarrow \mathbb{R}$, the Hamiltonian vector field $\vec{h}$ on $T^{*} M$ is defined by the identity: $d h=\sigma(\cdot, \vec{h})$.

Proposition 1. Let $\tilde{u} \in \mathcal{U}$ and $\tilde{q}(t)=q(t ; \tilde{u}), 0 \leq t \leq 1$; then $\tilde{u}$ is a normal extremal control if and only if there exists $\tilde{\lambda}_{t} \in T_{\tilde{q}(t)}^{*} M$ such that

$$
\begin{equation*}
\dot{\tilde{\lambda}}_{t}=\vec{h}_{\tilde{u}(t)}^{t},\left.\quad \frac{\partial h_{u}^{t}\left(\tilde{\lambda}_{t}\right)}{\partial u}\right|_{u=\tilde{u}(t)}=0, \quad 0 \leq t \leq 1 . \tag{4}
\end{equation*}
$$

This statement is standard, it is actually a weak version of the Pontryagin maximum principle.

A Lipschitzian curve $\tilde{\lambda}_{t}, 0 \leq t \leq 1$, in $T^{*} M$ that satisfies relations (4) is called a normal extremal associated to the normal extremal control $\tilde{u}$. The time-varying Hamiltonian system $\dot{\lambda}=\vec{h}_{\tilde{u}(t)}^{t}(\lambda)$ defines a flow

$$
\tilde{\Phi}^{t}: T^{*} M \rightarrow T^{*} M, \quad \tilde{\Phi}^{t}: \lambda(0) \mapsto \lambda(t),
$$

where $\dot{\lambda}(\tau)=\vec{h}_{\tilde{u}(\tau)}^{t}(\lambda(\tau)), 0 \leq \tau \leq t$. Obviously, $\tilde{\Phi}^{t}\left(\tilde{\lambda}_{0}\right)=\tilde{\lambda}_{t}$. Moreover, $\pi_{*} \vec{h}_{\tilde{u}(t)}^{t}(\lambda)=f_{\tilde{u}(t)}^{t}(q), \forall q \in M, \lambda \in T_{q}^{*} M$.

Let $\tilde{P}^{t}: M \rightarrow M$ be the flow generated by the time-varying system $\dot{q}=f_{\tilde{u}(t)}^{t}(q)$, i. e. $\tilde{P}^{t}: q(0) \mapsto q(t)$, where $\dot{q}(\tau)=f_{\tilde{u}(\tau)}^{\tau}(q(\tau)), 0 \leq \tau \leq t$. It follows that $\tilde{\Phi}^{t}$ are fiberwise transformations and $\tilde{\Phi}^{t}\left(T_{q}^{*} M\right)=T_{\tilde{P}^{t}(q)}^{*} M$. Moreover, the restriction of $\tilde{\Phi}^{t}$ to $T_{q}^{*} M$ is an affine map of the vector space $T_{q}^{*} M$ on the vector space $T_{P^{t}(q)}^{*} M$.

In what follows, we assume that $\tilde{u}$ is a normal extremal control and $\tilde{\lambda}_{t}$ is a normal extremal, they are fixed until the end of the paper. Then $\tilde{u}(t)$ is a critical point of the function $u \mapsto h_{u}^{t}\left(\tilde{\lambda}_{t}\right), u \in U$ for any $t \in[01]$. Hence the
 words, a self-adjoint linear map from $T_{\tilde{u}(t)} U$ to $T_{\tilde{u}(t)}^{*} U$.

Now we consider Hamiltonian functions

$$
g_{u}^{t}=\left(h_{u}^{t}-h_{\tilde{u}(t)}^{t}\right) \circ \tilde{\Phi}^{t}, \quad u \in U, t \in[0,1] .
$$

A time-varying Hamiltonian vector field $\bar{g}_{u}^{t}$ generates the flow $\left(\tilde{\Phi}^{t}\right)^{-1} \circ \Phi_{u}^{t}$, where $\Phi_{u}^{t}$ is the flow generated by the field $\vec{h}_{u}^{t}$. We have:

$$
g_{\tilde{u}(t)}^{t} \equiv 0, \quad \frac{\partial g_{\tilde{u}(t)}^{t}}{\partial u}\left(\tilde{\lambda}_{0}\right)=0, \quad \frac{\partial^{2} g_{\tilde{u}(t)}^{t}}{\partial^{2} u}\left(\tilde{\lambda}_{0}\right)=\frac{\partial^{2} h_{\tilde{u}(t)}^{t}}{\partial^{2} u}\left(\tilde{\lambda}_{t}\right) .
$$

We introduce a simplified notation $H_{t} \doteq \frac{\partial^{2} h_{\tilde{u}(t)}^{t}}{\partial^{2} u}\left(\tilde{\lambda}_{t}\right)$; recall that $H_{t}: T_{\tilde{u}(t)} U \rightarrow$ $T_{\tilde{u}(t)}^{*} U$ is a self-adjoint linear map.

More notations. Consider the map $u \mapsto \vec{g}_{u}^{t}\left(\tilde{\lambda}_{0}\right)$ from $U$ to $T_{\tilde{\lambda}_{0}}\left(T^{*} M\right) \mathrm{We}$ denote by $Z_{t}$ the differential of this map at $\tilde{u}(t), Z_{t} \doteq \frac{\partial \vec{g}_{\tilde{u}(t)}^{t}}{\partial u}\left(\tilde{\lambda}_{0}\right)$; then $Z_{t}$ is a linear map from $T_{\tilde{u}(t)} U$ to $T_{\tilde{\lambda}_{0}}\left(T^{*} M\right)$. We also set $X_{t}=\pi_{*} Z_{t}$; then $X_{t}$ is a linear map from $T_{\tilde{u}(t)} U$ to $T_{\tilde{q}(t)} M$. Finally, we denote by $J: T_{\tilde{\lambda}_{0}}\left(T^{*} M\right) \rightarrow$ $T_{\tilde{\lambda}_{0}}^{*}\left(T^{*} M\right)$ the anti-symmetric linear map defined by the identity $\sigma_{\tilde{\lambda}_{0}}(\cdot, \cdot)=$ $\langle J \cdot, \cdot\rangle$.

Note that $T_{u(\cdot)} \mathcal{U}$ is the space of measurable bounded mappings $t \mapsto$ $v(t) \in T_{u(t)} U, 0 \leq t \leq 1$. A simple rearrangement of the standard formula for the first variation in the optimal control theory (see for instance the textbook [1]) gives the identity:

$$
\left(D_{\tilde{u}} F_{t}\right) v=\tilde{P}_{*}^{t} \int_{0}^{t} X_{\tau} v(\tau) d \tau
$$

Assume that $\tilde{u}$ is a regular point of $F_{t}$ (i. e. $D_{\tilde{u}} F_{t}$ is surjective); then $\tilde{u}$ is a critical point of $\left.\varphi_{t}\right|_{F_{t}^{-1}(\tilde{q}(t))}$. We have:

$$
T_{\tilde{u}}\left(F_{t}^{-1}(\tilde{q}(t))=\operatorname{ker} D_{\tilde{u}} F_{t}=\left\{v \in T_{\tilde{u}} \mathcal{U}: \int_{0}^{t} X_{\tau} v(\tau) d \tau=0\right\}\right.
$$

The Hessian of $\left.\varphi_{1}\right|_{F_{1}^{-1}(\tilde{q}(1))}$ at $\tilde{u}$ is a quadratic form $D^{2} \varphi_{1}: \operatorname{ker} D_{\tilde{u}} F_{1} \rightarrow$ $\mathbb{R}$. This is the "second variation" of the optimal control problem at $\tilde{u}$, the main object of this paper. It has the following expression (see [1]):

$$
\begin{equation*}
D_{\tilde{u}}^{2} \varphi_{1}(v)=-\int_{0}^{1}\left\langle H_{t} v(t), v(t)\right\rangle d t-\int_{0}^{1}\left\langle J \int_{0}^{t} Z_{\tau} v(\tau) d \tau, Z_{t} v(t)\right\rangle d t \tag{5}
\end{equation*}
$$

where $v \in T_{\tilde{u}} \mathcal{U}$ and $\int_{0}^{1} X_{t} v(t) d t=0$.
Definition 2. The extremal $\tilde{\lambda}_{t}, 0 \leq t \leq 1$, is regular if $H_{t}$ is invertible for any $t \in[0,1]$ with a uniformly bounded inverse.

In what follows, we assume that the reference extremal is regular. Moreover, we assume that $H_{t}<0^{1}$ for any $t \in[0,1]$. This last assumption is motivated by the classical Legendre condition: for a regular extremal, integral quadratic form (5) has a finite negative inertia index if and only if $H_{t}<0$, for any $t \in[0,1]$.

Finally we introduce a "Gramm matrix", a self-adjoint linear map $\Gamma_{t}$ : $T_{q_{0}}^{*} M \rightarrow T_{q_{0}} M$ defined by the formula: $\Gamma_{t}=-\int_{0}^{t} X_{\tau} H_{t}^{-1} X_{\tau}^{*} d \tau$. We see that $\tilde{u}$ is a regular point of $F_{t}$ if and only if $\Gamma_{t}$ is invertible.

[^1]Remark 1. There is an apparently more general and natural way to define a control system. Let me briefly describe it confining within the timeinvariant case. Indeed, instead of the product $M \times U$ we may consider a locally trivial bundle over $M$ with a typical fiber $U$. Then $f$ is a smooth fiberwise map from this bundle to $T M$ that sends the fiber $U_{q}$ into $T_{q} M, q \in$ $M$. In this setting, the control and the correspondent trajectory is somehow a unique object: $u(t) \in U_{q(t)}$ and $\dot{q}(t)=f(u(t))$.

The situation is reduced to what we did before this remark if the bundle is trivial. In the general case, we simply trivialize the bundle in a neighborhood of the reference trajectory. To be more precise, we first take $\mathbb{R} \times M$, where $\mathbb{R}$ is the time axis and trivialize the pullback bundle over $\mathbb{R} \times M$ in a neighborhood of the graph of the reference trajectory (the graph does not have self-intersections while the original trajectory might have). Moreover, all statements of this paper use only $H_{t}, Z_{t}$, and $X_{t}, 0 \leq t \leq 1$, and these quantities depend only on the trivialization of the vector bundle $T_{\tilde{u}(t)} U_{q(t)}$ along (the graph of) the curve $u(t), 0 \leq t \leq 1$. A trivialization jf the vector bundle along such a curve is naturally achieved by a parallel transport in virtue of a linear connection on the vector bundle with the fibers $T_{u} U_{q}, q \in M, u \in U_{q}$ and the base $\bigcup_{q \in M} U_{q}$.

So, a linear connection on this vector bundle is actually all we need to write all the formulas. I leave to an interested reader to write them. I have decided not to follow this way in order to avoid an unnecessary language complication and to make the paper affordable for a larger audience.

## 3 Main results

According to our construction, the space $T_{\tilde{u}} \mathcal{U}$ consists of the $L^{\infty}$-maps $t \mapsto$ $v(t) \in T_{\tilde{u}(t)} U, 0 \leq t \leq 1$. At the same time, linear map $D_{\tilde{u}} F_{1}: T_{\tilde{u}} \mathcal{U} \rightarrow$ $T_{\tilde{q}(1)} M$ and quadratic form $D_{\tilde{u}}^{2} \varphi_{1}$ are continuous in a weaker topology $L^{2}$. Let $\mathcal{V}$ be the closure of $\operatorname{ker} D_{\tilde{u}} F_{1}$ in the topology $L^{2}$. Then $\mathcal{V}$ is a Hilbert space equipped with a Hilbert structure

$$
<v_{1} \mid v_{2}>\doteq \int_{0}^{1}\left\langle-H_{t} v_{1}(t), v_{2}(t)\right\rangle d t
$$

Formula (5) implies that

$$
\begin{equation*}
D_{\tilde{u}}^{2} \varphi_{1}(v)=<(I+K) v \mid v>, \quad v \in \mathcal{V} \tag{6}
\end{equation*}
$$

where $K$ is a compact symmetric operator on $\mathcal{V}$. In particular, the spectrum
of $K$ is real, the only limiting point of the spectrum is 0 , and any nonzero eigenvalue has a finite multiplicity.

What can we say about the "trace" of $K$ and the "determinant" of $I+K$ ? Let us consider a simple example, the least action principle for a charged particle in the plane in a constant magnetic field.

Example 1. Let $M=U=\mathbb{R}^{2}, q=\left(q^{1}, q^{2}\right), u=\left(u^{1}, u^{2}\right), f_{u}(q)=$ $u, \ell(q, u)=\frac{1}{2}|u|^{2}+r\left(q^{1} u^{2}-q^{2} u^{1}\right), q_{0}=0, \tilde{u}(t)=0$. Then

$$
\begin{gathered}
T^{*} M=\mathbb{R}^{2} \times \mathbb{R}^{2}=\left\{(p, q): p \in \mathbb{R}^{2}, q \in \mathbb{R}^{2}\right\}, \quad J(p, q)=(-q, p) \\
h_{u}(p, q)=\langle p, u\rangle-\frac{1}{2}|u|^{2}-r\left(q^{1} u^{2}-q^{2} u^{1}\right) \\
h_{\tilde{u}(t)}=0, \quad \tilde{\lambda}_{t}=(\tilde{p}(t), \tilde{q}(t))=(0,0), \quad g_{u}^{t}=h_{u}
\end{gathered}
$$

We have $H_{t} v=-v, Z_{t} v=\left(r v^{2},-r v^{1} ; v^{1}, v^{2}\right), X_{t} v=v, t \in[0,1], v \in \mathbb{R}^{2}$. Then

$$
\mathcal{V}=\left\{v \in L^{2}\left([0,1] ; \mathbb{R}^{2}\right): \int_{0}^{1} v(t) d t=0\right\}
$$

It is convenient to identify $\mathbb{R}^{2}$ with $\mathbb{C}$ as follows: $\left(v^{1}, v^{2}\right)=v^{1}+i v^{2}$. An immediate calculation gives the following expression for the operator $K$ :

$$
K v(t)=\int_{0}^{t} 2 \operatorname{riv}(\tau) d \tau-\int_{0}^{1} \int_{0}^{t} 2 \operatorname{riv}(\tau) d \tau d t
$$

The eigenfunctions of this operator have a form $t \mapsto c e^{2 \pi n i t}, 0 \leq t \leq 1, c \in$ $\mathbb{C}, n= \pm 1, \pm 2, \ldots$, where the eigenfunction $c e^{2 \pi n i t}$ corresponds to the eigenvalue $\frac{r}{\pi n}$.

We see that even in this model example the eigenvalues of $K$ do not form an absolutely convergent series and $K$ is not a trace class operator. On the other hand, the next theorem implies that a "principal value" of such a series does exist at least if the data are piece-wise real analytic.

Consider the operator $H_{t}^{-1} Z_{t}^{*} J Z_{t}: T_{\tilde{u}(t)} U \rightarrow T_{\tilde{u}(t)} U$. This operator is associated to an anti-symmetric bilinear form and has only purely imaginary eigenvalues. We denote by $\bar{\zeta}_{t}$ the sum of positive eigenvalues of the operator $i H_{t}^{-1} Z_{t}^{*} J Z_{t}$ counted according to multiplicity (here $i$ is the imaginary unit).

Remark 2. Analysing notations of Section 2 we see that $\bar{\zeta}_{t}$ is the sum of positive roots (counted according to multiplicity) of the equation

$$
\operatorname{det}\left(\left\{\frac{\partial h_{\tilde{u}(t)}^{t}}{\partial u}, \frac{\partial h_{\tilde{u}(t)}^{t}}{\partial u}\right\}\left(\tilde{\lambda}_{t}\right)+\operatorname{si} \frac{\partial^{2} h_{\tilde{u}(t)}^{t}}{\partial u^{2}}\left(\tilde{\lambda}_{t}\right)\right)=0
$$

with unknown $s$. Here $\{\cdot, \cdot\}$ is the Poisson bracket so that

$$
\left(v, v^{\prime}\right) \mapsto\left\{\frac{\partial h_{\tilde{u}(t)}^{t}}{\partial u} v, \frac{\partial h_{\tilde{u}(t)}^{t}}{\partial u} v^{\prime}\right\}\left(\tilde{\lambda}_{t}\right), \quad v, v^{\prime} \in T_{\tilde{u}(t)} U,
$$

is an anti-symmetric bilinear form and $\frac{\partial^{2} h_{\tilde{u}(t)}^{t}}{\partial u^{2}}\left(\tilde{\lambda}_{t}\right)$, is a symmetric bilinear form on $T_{\tilde{u}(t)} U$.

Let $\operatorname{Sp}(K) \subset \mathbb{R}$ be the spectrum of the operator $K, \operatorname{Sp}(K) \backslash\{0\}=$ $\mathrm{Sp}_{+}(K) \cup \mathrm{Sp}_{-}(K)$, where $\mathrm{Sp}_{ \pm}(K) \subset \mathbb{R}_{ \pm}$. Given $\alpha \in \operatorname{Sp}(K) \backslash 0$, we denote by $m_{\alpha}$ the multiplicity of the eigenvalue $\alpha$. Moreover, if $\mathrm{Sp}_{ \pm}(K)$ is an infinite set, then we introduce a natural ordering of $\operatorname{Sp}_{ \pm}(K)$ that is a monotone decreasing sequence $\alpha_{n}, n \in \mathbb{Z}_{ \pm}$, with the following properties:

$$
\begin{equation*}
\bigcup_{n \in \mathbb{Z}_{ \pm}}\left\{\alpha_{n}\right\}=\operatorname{Sp}_{ \pm}(K), \quad \#\left\{n \in \mathbb{Z}_{ \pm}: \alpha_{n}=\alpha\right\}=m_{\alpha} \tag{7}
\end{equation*}
$$

Theorem 1. If $\bar{\zeta}_{t} \equiv 0$, then $\alpha_{n}=O\left(|n|^{-2}\right)$ as $n \rightarrow \pm \infty$. If $\bar{\zeta}_{t}$ is not identical zero and $H_{t}$ and $Z_{t}$ are piecewise real analytic with respect to $t$ then $\mathrm{Sp}_{+}(K)$ and $\mathrm{Sp}_{-}(K)$ are both infinite and

$$
\begin{equation*}
\alpha_{n}=\frac{1}{\pi n} \int_{0}^{1} \bar{\zeta}_{t} d t+O\left(|n|^{-5 / 3}\right) \quad \text { as } \quad n \rightarrow \pm \infty . \tag{8}
\end{equation*}
$$

Remark 3. It is reasonable to expect that the statement of the theorem is valid without the piecewise-analyticity assumption. Moreover, the order $n^{-5 / 3}$ of the remainder term is certainly not optimal and perhaps can be substituted by $n^{-2}$. Anyway, the stated result is quite sufficient for our purposes while the proof of a stronger one would require a more sophisticated technique.

A cancellation of slow convergent to zero terms of the opposite sign in the expansion (8) gives the following:

Corollary 1. The depending on $\varepsilon>0$ families of real numbers

$$
\sum_{\substack{\alpha \in \operatorname{Sp}(K) \\|\alpha| \geq \varepsilon}} m_{\alpha} \alpha, \prod_{\substack{\alpha \in \operatorname{Sp}(K) \\|\alpha| \geq \varepsilon}}(1+\alpha)^{m_{\alpha}}
$$

have finite limits as $\varepsilon \rightarrow 0$.
We use natural notations for these limits:

$$
\operatorname{tr} K=\lim _{\varepsilon \rightarrow 0} \sum_{\substack{\alpha \in \operatorname{Sp}(K) \\|\alpha| \geq \varepsilon}} m_{\alpha} \alpha, \quad \operatorname{det}(I+K)=\lim _{\varepsilon \rightarrow 0} \prod_{\substack{\alpha \in \operatorname{Sp}(K) \\|\alpha| \geq \varepsilon}}(1+\alpha)^{m_{\alpha}} .
$$

We are going to compute these trace and determinant in terms of $H_{t}, Z_{t}, X_{t}$ and solutions of the following Jacobi system:

$$
\begin{equation*}
\dot{\eta}=-Z_{\tau} H_{\tau}^{-1} Z_{\tau}^{*} J \eta, \quad \eta(t) \in T_{\tilde{\lambda}_{0}}\left(T^{*} M\right), 0 \leq \tau \leq 1 . \tag{9}
\end{equation*}
$$

This is a linear time-varying Hamiltonian system in $T_{\tilde{\lambda}_{0}}\left(T^{*} M\right)$ associated to a nonnegative quadratic Hamiltonian $H_{t}(\eta)=-\frac{1}{2}\left\langle Z_{t}^{*} J \eta, H_{t}^{-1} Z_{t}^{*} J \eta\right\rangle$.

In what follows, we identify the space $T_{q_{0}}^{*} M$ with its tangent $T_{\tilde{\lambda}_{0}}\left(T_{q_{0}}^{*} M\right) \subset$ $T_{\tilde{\lambda}_{0}}\left(T^{*} M\right)$ and introduce linear maps $Q_{t}: T_{q_{0}}^{*} M \rightarrow T_{q_{0}} M, 0 \leq t \leq 1$, by the formula: $Q_{t}(\eta(0))=\pi_{*} \eta(t)$, where $\eta(\tau), 0 \leq \tau \leq t$, is a solution of equation (9) and $\eta(0) \in T_{q_{0}}^{*} M$.

Remark 4. According to our assumptions, $\tilde{u}(\tau)$ is a strict local maximum of the function $u \mapsto h_{u}^{\tau}\left(\tilde{\lambda}_{\tau}\right)$. Assume that this is a global maximum and moreover $h^{\tau}(\lambda)=\max _{u \in U} h_{u}^{\tau}(\lambda)$ is smooth with respect to $\lambda \in T^{*} M$. Then $\tilde{\lambda}_{\tau}, 0 \leq \tau \leq 1$, is a solution of the Hamiltonian system $\dot{\lambda}=\vec{h}^{\tau}(\lambda)$. We define the exponential map $\mathcal{E}_{q}^{t}: T_{q}^{*} M \rightarrow M$ by the formula $\mathcal{E}_{q}^{t}\left(\lambda_{0}\right)=\pi\left(\lambda_{t}\right)$, where $\dot{\lambda}_{\tau}=\vec{h}^{\tau}\left(\lambda_{\tau}\right), 0 \leq \tau \leq t$. It is not hard to show that $Q_{t}=\left(\tilde{P}_{*}^{t}\right)^{-1} D_{\tilde{\lambda}_{0}} \mathcal{E}_{q_{0}}^{t}$.

Theorem 2. Under conditions of Theorem 1, the following identities are valid:

$$
\begin{gathered}
\operatorname{det}(I+K)=\operatorname{det}\left(Q_{1} \Gamma_{1}^{-1}\right) \\
\operatorname{tr} K=\operatorname{tr}\left(\int_{0}^{1} \int_{0}^{t} X_{t} H_{t}^{-1} Z_{t}^{*} J Z_{\tau} H_{\tau}^{-1} X_{\tau}^{*} d \tau d t \Gamma_{1}^{-1}\right) .
\end{gathered}
$$

Let us apply this theorem to Example 1. In our coordinates, $\Gamma_{1}=I$. We have to find matrix $Q_{1}$. It is convenient to use complex notations: $\eta=(p ; q) \in \mathbb{C} \times \mathbb{C}, p=i p_{1}+p_{2}, q=i q_{1}+q_{2}$. System (8) has the form: $\dot{p}=\operatorname{irp}-\nu^{2} q, \quad \dot{q}=i r q+p$. We have to find $q(t)$ under conditions $p(0)=1, q(0)=0$; then $Q_{1}$ is just complex number $q(1)$ treated as a $2 \times 2$ real matrix, $\operatorname{det} Q_{1}=|q(1)|^{2}$. A simple calculation gives: $q(t)=\frac{\sin \nu t}{\nu} e^{i \nu t}$. Keeping in mind that all eigenvalues of the operator $K$ have multiplicity 2, we obtain the square of a classical Euler identity: $\prod_{n=1}^{+\infty}\left(1-\left(\frac{r}{\pi n}\right)^{2}\right)^{2}=$ $\left(\frac{\sin r}{r}\right)^{2}$.

Now we consider one more very simple example, a harmonic oscillator.
Example 2. $M=U=\mathbb{R}, f_{u}(q)=u, \ell(q, u)=\frac{1}{2}\left(u^{2}-r q^{2}\right), q_{0}=$ $0, \tilde{u}(t)=0$. Then $h_{u}(p, q)=p u-\frac{1}{2}\left(u^{2}-r q^{2}\right), h_{\tilde{u}(t)}=\frac{r}{2} q^{2}$,
$\tilde{\Phi}_{t}(p, q)=\binom{p-t r q}{q}, g_{u}^{t}=(p-t r q) u-\frac{1}{2} u^{2}, \quad H_{t}=-1, Z_{t}=\binom{t r}{1}, X_{t}=1$.

Operator $K$ has a form:

$$
K v(t)=r \int_{0}^{t}(t-\tau) v(\tau) d \tau-r \int_{0}^{1} \int_{0}^{t}(t-\tau) v(\tau) d \tau d t
$$

The eigenfunctions of this operator have a form $t \mapsto c \cos (\pi n t), c \in \mathbb{R}, n=$ $1,2, \ldots$, where the eigenfunction $c \cos (\pi n t)$ corresponds to the eigenvalue $-\frac{r}{(\pi n)^{2}}$. Moreover, $Q_{1}=\frac{\sin \sqrt{r}}{\sqrt{r}}$ if $r>0$ and $Q_{1}=\frac{\operatorname{sh} \sqrt{|r|}}{\sqrt{|r|}}$ if $r<0$. The determinant formula from Theorem 2 coincides with the Euler identity: $\prod_{n=1}^{\infty}\left(1-\frac{r}{(\pi n)^{2}}\right)=\frac{\sin \sqrt{r}}{\sqrt{r}}$ or its hyperbolic version. The trace formula gives another famous Euler observation: $\sum_{n=1}^{\infty} \frac{r}{(\pi n)^{2}}=\frac{r}{6}$.

Harmonic oscillator is a special case of a more general example where the spectrum is still explicitly computed.

Example 3. $M=U=\mathbb{R}^{m}, f_{u}(q)=A q+u, \quad \ell(q, u)=\frac{1}{2}\left(|u|^{2}-\right.$ $\langle q, R q\rangle), q_{0}=0, \tilde{u}(t)=0$, where $A$ is a $m \times m$-matrix and $R$ is a symmetric $m \times m$-matrix. Then $h_{u}(p, q)=\langle p, A q+u\rangle-\frac{1}{2}\left(|u|^{2}-\langle q, R q\rangle\right), \quad h_{\tilde{u}(t)}=$ $\langle p, A q\rangle+\frac{1}{2}\langle q, R q\rangle$,

$$
\begin{gathered}
\tilde{\Phi}_{t}(p, q)=\binom{e^{-t A^{*}} p-\int_{0}^{t} e^{(\tau-t) A^{*}} R e^{\tau A} d \tau}{e^{t A} q}, \\
g_{u}^{t}=\left\langle e^{-t A^{*}} p-\int_{0}^{t} e^{(\tau-t) A^{*}} R e^{\tau A} d \tau, u\right\rangle-\frac{1}{2}|u|^{2}, \\
H_{t}=-I, \quad Z_{t}=\binom{\int_{0}^{t} e^{\tau A^{*}} R e^{\tau A} d \tau}{e^{-t A}}, \quad X_{t}=e^{-t A} .
\end{gathered}
$$

Then $\mathcal{V}=\left\{v \in L^{2}\left([0,1] ; \mathbb{R}^{2}\right): \int_{0}^{1} e^{-t A} v(t) d t=0\right\}$,

$$
<K v \mid v>=\int_{0}^{1}\left\langle\iint_{\Delta \leq \tau_{1} \leq \tau_{2} \leq t} e^{\left(\tau_{2}-t\right) A^{*}} R e^{\left(\tau_{2}-\tau_{1}\right) A} v\left(\tau_{1}\right) d \tau_{2} d \tau_{1}, v(t)\right\rangle d t
$$

A vector function $v(\cdot) \in \mathcal{V}$ is an eigenfunction of the operator $K$ with an eigenvalue $\alpha \in \mathbb{R}$ if and only if

$$
\begin{equation*}
\alpha v(t)=\iint_{0 \leq \tau_{1} \leq \tau_{2} \leq t} e^{\left(\tau_{2}-t\right) A^{*}} R e^{\left(\tau_{2}-\tau_{1}\right) A} v\left(\tau_{1}\right) d \tau_{2} d \tau_{1}+e^{-t A^{*}} c, \quad t \in[0,1], \tag{10}
\end{equation*}
$$

for some constant vector $c \in \mathbb{R}^{m}$. We denote: $y(t)=\int_{0}^{t} e^{(t-\tau) A} v(\tau) d \tau$; then we differentiate twice equation (10) and obtain that this equation is equivalent to the boundary values problem $y(0)=y(1)=0$ for the ordinary differential equation

$$
\alpha \ddot{y}=\alpha\left(A-A^{*}\right) \dot{y}+\left(\alpha A^{*} A+R\right) y .
$$

From now on we assume that the matrix $A$ is symmetric, $A^{*}=A$; then nonzero eigenvalues and corresponding eigenfunctions are nontrivial solution of the boundary value problem:

$$
\ddot{y}=\left(A^{2}+\frac{1}{\alpha} R\right) y, \quad y(0)=y(1)=0 .
$$

We obtain that $\alpha$ is a nonzero eigenvalue of $K$ if and only if there exists a positive integer $n$ such that $-(\pi n)^{2}$ is an eigenvalue of the matrix $A^{2}+\frac{1}{\alpha} R$. Moreover, the multiplicities of the eigenvalue $\alpha$ of $K$ and of the eigenvalue $-(\pi n)^{2}$ of $A^{2}+\frac{1}{\alpha} R$ are equal. In other words, $\operatorname{det} K=\prod_{n=1}^{\infty} \prod_{i}\left(1-s_{n i}\right)$ where $s_{n i}$ are real roots (counted with their multiplicities) of the following polynomial equation with unknown $s$ :

$$
\operatorname{det}\left(s\left(A^{2}+(\pi n)^{2} I\right)-R\right)=0
$$

Actually, all roots of this polynomial are real since $A^{2}+(\pi n)^{2} I$ is a signdefinite symmetric matrix. Let $\phi(s)=\operatorname{det}\left(s\left(A^{2}+(\pi n)^{2} I\right)-R\right), \phi(s)=$ $\sum_{i=1}^{m} c_{n i} s^{i}$; then $c_{n m}=\operatorname{det}\left(A^{2}+(\pi n)^{2} I\right)$ and $\prod_{n=1}^{\infty} \prod_{i}\left(1-s_{n i}\right)=\frac{\phi(1)}{c_{n m}}$. We obtain that

$$
\prod_{i=1}^{m}\left(1-s_{n i}\right)=\frac{1}{c_{n m}} \operatorname{det}\left(\left(A^{2}+(\pi n)^{2} I\right)-R\right)=\operatorname{det}\left(I-R\left(A^{2}+(\pi n)^{2} I\right)^{-1}\right)
$$

Hence $\operatorname{det} K=\prod_{n=1}^{\infty} \operatorname{det}\left(I-R\left(A^{2}+(\pi n)^{2} I\right)^{-1}\right)$.
Similarly, $\operatorname{tr} K=-\sum_{n=1}^{\infty} \operatorname{tr}\left(R\left(A^{2}+(\pi n)^{2} I\right)^{-1}\right)$.
Now we have to compute the right-hand sides of the identities of Theorem 2. First of all, $\Gamma_{1}=\int_{0}^{1} e^{-2 t A} d t$. Then we consider Jacobi system (9). We have $T^{*} M=\mathbb{R}^{m} \times \mathbb{R}^{m}=\left\{(p, q): p, q \in \mathbb{R}^{m}\right\}$; then $Q_{t} p_{0}=q_{t}$ where $t \mapsto\left(p_{t}, q_{t}\right)$ is the solution of the Jacobi system with the initial value $q_{0}=0$. We set $y(t)=e^{t A} q_{t}$, differentiate in virtue of the Jacobi system and obtain the equation: $\ddot{y}=\left(A^{2}-R\right) y$. Moreover, $\dot{y}(0)=p_{0}$. Hence
$Q_{t}=\left(R-A^{2}\right)^{-1 / 2} \sin \left(\left(A^{2}-R\right)^{1 / 2} t\right)$. This formula is valid also for a signindefinite matrix $R-A^{2}$ if we properly interpret the square root or simply make the computations in coordinates where the matrix $A^{2}-R$ is diagonal. Finally, putting together all the formulas, we obtain the following generalization of the Euler identities:

Proposition 2. Let $A, R$ be symmetric matrices. Then:

$$
\begin{gathered}
\prod_{n=1}^{\infty} \operatorname{det}\left(I-R\left(A^{2}+(\pi n)^{2} I\right)^{-1}\right)=\frac{2 \operatorname{det}\left(\sin \sqrt{R-A^{2}}\right)}{\operatorname{det}\left(\sqrt{R-A^{2}} \int_{-1}^{1} e^{t A} d t\right)}, \\
\sum_{n=1}^{\infty} \operatorname{tr}\left(R\left(A^{2}+(\pi n)^{2} I\right)^{-1}\right)= \\
\operatorname{tr}\left(\iiint_{0 \leq \tau_{1} \leq \tau_{2} \leq t \leq 1} e^{\left(\tau_{2}-2 t\right) A} R e^{\left(\tau_{2}-2 \tau_{1}\right)} d \tau_{2} d \tau_{1} d t\left(\int_{0}^{1} e^{-2 t A} d t\right)^{-1}\right) .
\end{gathered}
$$

The right-hand side of the determinant formula has an obvious meaning also in the case of a degenerate $R-A^{2}$. If $m=1, A=a, R=a^{2}+b^{2}$, we get:

$$
\prod_{n=1}^{\infty}\left(1-\frac{a^{2}+b^{2}}{a^{2}+(\pi n)^{2}}\right)=\frac{a \sin b}{b \operatorname{sh} a}
$$

an interpolation between the classical Euler identity and its hyperbolic version. The trace identity is essentially simplified if the matrices $R$ and $A$ commute. In the commutative case we obtain:

$$
\sum_{n=1}^{\infty} \operatorname{tr}\left(R\left(A^{2}+(\pi n)^{2} I\right)^{-1}\right)=\frac{1}{2} \operatorname{tr}\left(R(A \operatorname{cth} A-I) A^{-2}\right)
$$

## 4 Proof of Theorem 1

We start with some definitions and notations. A compact quadratic form $b$ on the Hilbert space $\mathcal{V}$ is a form defined by a compact symmetric operator $B, b(v)=<B v \mid v>, v \in \mathcal{V}$. The spectrum of $b$ and its positive and negative parts are, by the definition, those of $B$, i. e. $\mathrm{Sp}_{ \pm}(b)=\mathrm{Sp}_{ \pm}(B)$. Recall that the eigenvalues of $B$ are just critical values of the restriction of $b$ to the unit sphere in $\mathcal{V}$.

Let $\mathcal{V}_{0} \subset \mathcal{V}$ be a Hilbert subspace, then the form $\left.b\right|_{\mathcal{V}_{0}}$ is defined by the composition of $\left.B\right|_{\mathcal{V}_{0}}$ and the orthogonal projection of $\mathcal{V}$ on $\mathcal{V}_{0}$. The relation between $\operatorname{Sp}(b)$ and $\operatorname{Sp}\left(b \mid \nu_{0}\right)$ is ruled by the classical Rayleigh-Courant minimax principle (see $[4,5]$ ).

Assume that $\mathrm{Sp}_{ \pm}(b)$ is infinite; then a natural ordering $\beta_{n}, n \in \mathbb{Z}_{ \pm}$, of $\mathrm{Sp}_{ \pm}(b)$ is defined in the same way as the natural ordering of $\mathrm{Sp}_{ \pm}(K)$ (see (7)).

Definition 3. We say that b has the spectrum of capacity $\varsigma>0$ with the remainder of order $\nu>1$ if $\mathrm{Sp}_{+}(b)$ and $\mathrm{Sp}_{-}(b)$ are both infinite and

$$
\begin{equation*}
\beta_{n}=\frac{\varsigma}{n}+O\left(n^{-\nu}\right) \quad \text { as } \quad n \rightarrow \pm \infty \tag{11}
\end{equation*}
$$

We say that b has the spectrum of zero capacity with the remainder of order $\nu$ if either $\mathrm{Sp}_{ \pm}(b)$ is finite or $\beta_{n}=O\left(n^{-\nu}\right)$ as $n \rightarrow \pm \infty$.

Let $b_{i}$ be a quadratic form on the Hilbert space $\mathcal{V}_{i}, i=1,2$; then $b_{1} \oplus b_{2}$ is a quadratic form on $\mathcal{V}_{1} \oplus \mathcal{V}_{2}, \operatorname{Sp}\left(b_{1} \oplus b_{2}\right)=\operatorname{Sp}\left(b_{1}\right) \cup \operatorname{Sp}\left(b_{2}\right)$ and the multiplicities of common eigenvalues are added.
Proposition 3. (i) If $b$ has the spectrum of capacity $\varsigma \geq 0$, then sb has the spectrum of capacity ss with the remainder of the same order as b, for any $s \in \mathbb{R}$.
(ii) If $b_{1}, b_{2}$ have the spectra of capacities $\varsigma_{1}, \varsigma_{2}$ with the remainders of equal orders, then $b_{1} \oplus b_{2}$ has the spectrum of capacity $\varsigma_{1}+\varsigma_{2}$ and the remainder of the same order as $b_{1}, b_{2}$.
(iii) Let $\mathcal{V}_{0}$ be a Hilbert subspace of the Hilbert space $\mathcal{V}$ and $\operatorname{dim}\left(\mathcal{V} / \mathcal{V}_{0}\right)<\infty$. Assume that one of two forms $b$ or $b \mid \mathcal{V}_{0}$ has the spectrum of capacity $\varsigma \geq 0$ with a remainder of order $\nu \leq 2$. Then the second form has the spectrum of the the same capacity $\varsigma$ with a remainder of the same order $\nu$.
(iv) Let the forms $b$ and $\hat{b}$ be defined on the same Hilbert space $\mathcal{V}$, where $b$ has the spectrum of capacity $\varsigma$ and $\hat{b}$ has the spectrum of zero capacity, both with the reminder term of order $\nu$. Then the form $b+\hat{b}$ has the spectrum of capacity $\varsigma$ with the reminder term of order $\frac{2 \nu+1}{\nu+1}$.
Proof. Statement (i) is obvious. To prove (ii) we re-write asymptotic relation (11) in a more convenient form. An equivalent relation for positive $n$ reads:

$$
\#\left\{k \in \mathbb{Z}: 0<\frac{1}{\beta_{k}}<n\right\}=\varsigma n+O\left(n^{2-\nu}\right), \quad \text { as } \quad n \rightarrow \infty
$$

and similarly for negative $n$. Statement (ii) follows immediately.
Statement (iii) follows from the Rayleigh-Courant minimax principle for the eigenvalues and the relation: $\left|\frac{\varsigma}{n}-\frac{\varsigma}{n+j}\right|=O\left(\frac{1}{n^{2}}\right)$ as $|n| \rightarrow \infty$ for any fixed $j$.

To prove (iv) we use the Weyl inequality for the eigenvalues of the sum of two forms. Weyl inequality is a straightforward corollary of the minimax principle, it claims that the positive eigenvalue number $i+j-1$ in the natural ordering of the sum of two forms does not exceed the sum of the eigenvalue number $i$ of the first summand and the eigenvalue number $j$ of the second summand. Of course, we may equally works with naturally ordered negative eigenvalues simply changing the signs of the forms.

In our case, to have both sides estimates we first present $b+\hat{b}$ as the sum of $b$ and $\hat{b}$ and then present $b$ as the sum of $b+\hat{b}$ and $-\hat{b}$. In the first case we apply the Weyl inequality with $i=n-\left[n^{\delta}\right], j=\left[n^{\delta}\right]$ for some $\delta \in(0,1)$, and in the second case we take $i=n, j=\left[n^{\delta}\right]$. The best result is obtained for $\delta=\frac{1}{\nu+1}$.

We have to prove that the spectrum of operator $K$ (see (6), (5)) has capacity $\frac{1}{\pi} \int_{0}^{1} \bar{\zeta}_{t} d t$ with the remainder of order $\frac{5}{3}$. First, we may identify $T_{\tilde{u}(t)}$ with $\mathbb{R}^{m}$ and assume that $H_{t}=-I$. Indeed, if we trivialize the vector bundle $\tilde{u}^{*}(T U)$ over the segment $[0,1]$, then $H_{t}$ becomes a negative definite symmetric matrix. Then we substitute $v$ by $\left(-H_{t}\right)^{\frac{1}{2}} v$ and $Z_{t}$ by $Z_{t}\left(-H_{t}\right)^{-\frac{1}{2}}$.

Let $0=t_{0}<t_{1}<\cdots<t_{l}<t_{l+1}=1$ be a subdivision of the segment $[0,1]$. The subspace

$$
\left\{v \in \mathcal{V}: \int_{t_{i}}^{t_{i+1}} X_{t} v(t) d t=0, i=0,1, \ldots, l\right\}=\bigoplus_{i=0}^{l} \mathcal{V}_{i}, \mathcal{V}_{i} \subset L^{2}\left([0,1] ; \mathbb{R}^{m}\right)
$$

has a finite codimension in $\mathcal{V}$. The quadratic form

$$
\begin{equation*}
\langle K v \mid v\rangle=\int_{0}^{1}\left\langle J Z_{t} v(t), \int_{0}^{t} Z_{\tau} v(\tau) d \tau\right\rangle d t \tag{12}
\end{equation*}
$$

restricted to this subspace turns into the direct sum of the forms

$$
<K_{i} v \mid v>=\int_{t_{i}}^{t_{i+1}}\left\langle J Z_{t} v(t), \int_{t_{i}}^{t} Z_{\tau} v(\tau) d \tau\right\rangle d t, \quad v_{i} \in \mathcal{V}_{i}, i=0,1 \ldots l .
$$

Indeed, the relations $\int_{t_{i}}^{t_{i+1}} X_{t} v(t) d t=0$ imply that

$$
\left\langle\int_{t_{i}}^{t_{i+1}} J Z_{t} v(t) d t, \int_{t_{j}}^{t_{j+1}} Z_{t} v(t) d t\right\rangle=0 .
$$

According to Proposition 3 it is sufficient to prove our theorem for the operators $K_{i}, i=0,1 \ldots, l$. In particular, we may substitute the piecewise analyticity assumption in the statement of the theorem by the analyticity one.

Moreover, under the analyticity condition we may assume that

$$
Z_{t}^{*} J Z_{t}=\bigoplus_{j=1}^{k}\left(\begin{array}{cc}
0 & -\zeta_{j}(t) \\
\zeta_{j}(t) & 0
\end{array}\right)
$$

where $0 \leq 2 k \leq m$ and $\zeta_{j}(t)$ are not identical zero. Indeed, according to the Rayleigh theorem (see [5]), there exists an analytically depending on $t$ orthonormal basis in which our anti-symmetric matrix takes a desired form.

The functions $\zeta_{j}(t), j=1, \ldots, k$, are analytic and may have only isolated zeros. Hence we may take a finer subdivision of $[0,1]$ in such a way that $\zeta_{j}(t), j=1, \ldots, k$, do not change sign on the segments $\left[t_{i}, t_{i+1}\right]$. Moreover, a simple change of the basis of $\mathbb{R}^{m}$ if necessary allows us to assume that $\zeta_{j}(t) \geq 0, t \in\left[t_{i}, t_{i+1}\right]$. Actually, to simplify notations a little bit, we may simply assume that $\zeta_{j}(t) \geq 0,0 \leq t \leq 1, j=1, \ldots, k$. In this case $\bar{\zeta}(t)=\sum_{j=1}^{k} \zeta_{j}(t)$.

Let us study quadratic form (12) on the space $\left\{v \in L^{2}\left([0,1] ; \mathrm{R}^{m}\right)\right.$ : $\left.\int_{0}^{1} v(t) d t=0\right\}$. Recall that we are allowed by Proposition 3 to work on any subspace of $L^{2}\left([0,1] ; \mathrm{R}^{m}\right)$ of a finite codimension. We set $w(t)=\int_{0}^{t} v(\tau) d \tau$; a double integration by parts gives:

$$
\begin{aligned}
& \int_{0}^{1}\left\langle J Z_{t} v(t), \int_{0}^{t} Z_{\tau} v(\tau) d \tau\right\rangle d t=\int_{0}^{1}\left\langle J Z_{t} v(t), Z_{t} w(t)\right\rangle d t+ \\
& \int_{0}^{1}\left\langle J Z_{t} w(t), \dot{Z}_{t} w(t)\right\rangle d t+\int_{0}^{1}\left\langle J \dot{Z}_{t} w(t), \int_{0}^{t} \dot{Z}_{\tau} w(\tau) d \tau\right\rangle d t
\end{aligned}
$$

Moreover, we have:

$$
\left|\int_{0}^{1}\left\langle J Z_{t} w(t), \dot{Z}_{t} w(t)\right\rangle d t+\int_{0}^{1}\left\langle J \dot{Z}_{t} w(t), \int_{0}^{t} \dot{Z}_{\tau} w(\tau) d \tau\right\rangle d t\right| \leq c \int_{0}^{1}|w(t)|^{2} d t
$$

for some constant $c$.
Let $\bar{\lambda}_{n}$ and $\underline{\lambda}_{n}, n \in \mathbb{Z} \backslash\{0\}$, be naturally ordered non zero eigenvalues of the quadratic forms

$$
\int_{0}^{1}\left\langle J Z_{t} v(t), Z_{t} w(t)\right\rangle d t+c \int_{0}^{1}|w(t)|^{2} d t
$$

and

$$
\int_{0}^{1}\left\langle J Z_{t} v(t), Z_{t} w(t)\right\rangle d t-c \int_{0}^{1}|w(t)|^{2} d t
$$

correspondently. The minimax principle for the eigenvalues implies that $\underline{\lambda}_{n} \leq \alpha_{n} \leq \bar{\lambda}_{n}, n \in \mathbb{Z} \backslash\{0\}$. Moreover, the form

$$
\int_{0}^{1}\left\langle J Z_{t} v(t), Z_{t} w(t)\right\rangle d t \pm c \int_{0}^{1}|w(t)|^{2} d t
$$

splits in the direct sum of the forms

$$
\begin{equation*}
\int_{0}^{1} \zeta_{j}(t)\left\langle J v_{j}(t), w_{j}(t)\right\rangle d t \pm c \int_{0}^{1}\left|w_{j}(t)\right|^{2} d t, \quad j=1, \ldots, k \tag{13}
\end{equation*}
$$

where $v_{j}(t) \in \mathbb{R}^{2}, w_{j}(t)=\int_{0}^{t} v_{j}(t), w_{j}(1)=0, J=\left(\begin{array}{cc}0 & -1 \\ 1 & 0\end{array}\right)$ and we simply ignore the identically vanishing part if $2 k<m$.

It remains to estimate the spectrum of the forms (13). To do that, we study the spectrum of the forms

$$
\begin{equation*}
\int_{0}^{1} \zeta_{j}(t)\left\langle J v_{j}(t), w_{j}(t)\right\rangle d t, \quad \int_{0}^{1}\left|w_{j}(t)\right|^{2} d t \tag{14}
\end{equation*}
$$

and apply the statement (iv) of Proposition 3 to their linear combinations. To simplify calculations, we identify $\mathbb{R}^{2}$ with $\mathbb{C}$ as follows: $\left(v_{j}^{1}, v_{j}^{2}\right)=v_{j}^{1}+i v_{j}^{2}$, then $J$ is the multiplication on the imaginary unit $i$. A complex-valued function $v_{j}(\cdot)$ is an eigenfunction with an eigenvalue $\lambda$ for the first of two forms (14) if and only if it is a critical point of the functional

$$
\int_{0}^{1}\left\langle\zeta(t) i v_{j}(t), w_{j}(t)\right\rangle-\lambda\left\langle v_{j}(t), v_{j}(t)\right\rangle d t \quad w_{j}(0)=w_{j}(1)=0
$$

The Euler-Lagrange equation for this functional reads:

$$
2 \lambda \ddot{w}_{j}+\left(\zeta_{j} i w_{j}\right)^{\cdot}+\zeta_{j} i \dot{w}=0, \quad w_{j}(0)=w_{j}(1)=0 .
$$

We set $x(t)=e^{\frac{i}{2 \lambda} \int_{0}^{t} \zeta_{j}(\tau) d \tau} w_{j}(t)$, plug-in this expression in the equation and arrive to a standard Sturm-Liouville problem:

$$
\ddot{x}+\left(\frac{\zeta_{j}}{2 \lambda}\right)^{2} x=0, \quad x(0)=x(1)=0
$$

We see that the spectrum is double (recall that $x \in \mathbb{C}$ ) and is symmetric with respect to the origin. The asymptotics of the spectrum for the SturmLiouville problem is well-known (see [4]). We obtain:

$$
\lambda_{n}=\frac{1}{\pi n} \int_{0}^{1} \zeta_{j}(t) d t+O\left(\frac{1}{n^{2}}\right), \quad n \rightarrow \pm \infty
$$

In other words, this spectrum has capacity $\frac{1}{\pi} \int_{0}^{1} \zeta_{j}(t) d t$ and the remainder term of order 2. The Euler-Lagrange equation for the eigenfunctions problem for the second of two forms (14) reads: $\lambda \ddot{w}_{j}+w_{j}=0, w_{j}(0)=w_{j}(1)=0$; hence $\lambda_{n}=\frac{1}{(\pi n)^{2}}=O\left(\frac{1}{n^{2}}\right)$. The spectrum of this form has zero capacity with the remainder of order 2 .

## 5 Proof of Theorem 2

We again assume that $H_{t}=-I$ using the same preliminary change of variables as in the proof of Theorem 1, if necessary: substitute $v$ by $(-H)^{\frac{1}{2}} v$ and $Z_{t}$ by $Z_{t}(-H)^{-\frac{1}{2}}$. Moreover, we fix some coordinates in a neighborhood of $q_{0} \in M$ and use induced coordinates in $T^{*} M$ so that $T_{\tilde{\lambda}_{0}}\left(T^{*} M\right)$ is identified with $\mathbb{R}^{d} \times \mathbb{R}^{d}=\left\{(p, q): p, q \in \mathbb{R}^{d}\right\}$, where $d=\operatorname{dim} M, \pi_{*}(p, q)=q$, and

$$
\sigma_{\tilde{\lambda}_{0}}\left((p, q),\left(p^{\prime}, q^{\prime}\right)\right)=\left\langle J\binom{p}{q},\binom{p^{\prime}}{q^{\prime}}\right\rangle, \quad J=\left(\begin{array}{cc}
0 & -I \\
I & 0
\end{array}\right) .
$$

The map $Z_{t}: \mathbb{R}^{m} \rightarrow \mathbb{R}^{d} \times \mathbb{R}^{d}$ has a form $Z_{t}: v \mapsto\left(Y_{t} v, X_{t} v\right)$, where $Y_{t}$ and $X_{t}$ are $d \times m$-matrices. Let $s \in \mathbb{C}$; we define $Z_{t}^{s}: \mathbb{C}^{d} \times C^{d}$ by the formula $Z_{t}^{s} v=\left(s Y_{t} v, X_{t} v\right)$. Now consider the complexified Jacobi equation:

$$
\begin{equation*}
\dot{\eta}=Z_{t}^{s} Z_{t}^{s *} J \eta, \quad \eta \in \mathbb{C}^{d} \times \mathbb{C}^{d}, \tag{15}
\end{equation*}
$$

where the transposition "*" corresponds to the complex inner product and not to the Hermitian one! The matrix form of this equation is as follows:

$$
\begin{equation*}
\binom{\dot{p}}{\dot{q}}=\binom{s Y_{t}}{X_{t}}\left(s Y_{t}^{*}, X_{t}^{*}\right)\binom{-q}{p} . \tag{16}
\end{equation*}
$$

Proposition 4. Let $\Phi_{t}^{s}: \mathbb{C}^{d} \times \mathbb{C}^{d} \rightarrow \mathbb{C}^{d} \times \mathbb{C}^{d}, 0 \leq t \leq 1$ be the fundamental matrix of the complexified Jacobi equation, i.e.

$$
\frac{d}{d t} \Phi_{t}^{s}=Z_{t}^{s} Z_{t}^{s *} J \Phi_{t}^{s}, \quad \Phi_{0}^{s}=I
$$

then $s \mapsto \Phi_{1}^{s}, s \in \mathbb{C}$, is an entire matrix-valued function and

$$
\left\|\Phi_{1}^{s}\right\| \leq c e^{c|s|}
$$

for some constant c
Proof. The Volterra series for our system has the form:

$$
\Phi_{1}^{s}=I+\sum_{n=1}^{\infty} \int_{0}^{1} \int_{0}^{t_{1}} \cdots \int_{0}^{t_{n}-1} Z_{t_{1}}^{s} Z_{t_{1}}^{s *} J \cdots Z_{t_{n}}^{s} Z_{t_{n}}^{s *} J d t_{n} \ldots d t_{1} .
$$

The $n$-th term of the series is a polynomial matrix of $s$ whose norm is bounded by $\frac{1}{n!}\left(\int_{0}^{1}\left\|Z_{t}^{s} Z_{t}^{s *}\right\| d t\right)^{n}$ for any $s$ and $n$. The sum of the series is an entire function by the Weierstrass theorem. Moreover, we obtain the estimate: $\left\|\Phi_{t}^{s}\right\| \leq e^{\int_{0}^{1}\left\|Z_{t}^{s} Z_{t}^{s *}\right\| d t}$ but it is worse than one we need.

To obtain a better estimate we make a change of variables in the matrix form of the complexified Jacobi equation (see (16)). We set $\eta^{s}=(p, s q)$; then $\dot{\eta}^{s}=s Z_{t} Z_{t}^{*} J \eta^{s}$ and we obtain: $\left|\eta^{s}(1)\right| \leq e^{|s|} \int_{0}^{1}\left\|Z_{t} Z_{t}^{*}\right\| d t\left|\eta^{s}(0)\right|$ that gives us the desired estimate of $\left\|\Phi_{1}^{s}\right\|$ for $s$ separated from zero.

We define linear maps $Q_{t}^{s}: \mathbb{C}^{n} \rightarrow \mathbb{C}^{n}$ by the formula: $Q_{t}^{s} p(0)=q(t)$, where $t \mapsto\binom{p(t)}{q(t)}$ is the solution of system (16) with the initial condition $q(0)=0$. Note that the matrix $Q_{t}^{s}$ is real if $s \in \mathbb{R}$; moreover, $\left.Q_{t}^{1}\right|_{\mathbb{R}^{d}}=$ $Q_{t}$ (linear maps $Q_{t}$ were defined just before the statement of Theorem 2). Proposition 4 implies that $\left\|Q_{1}^{s}\right\| \leq c e^{c|s|}$.

Proposition 5. Let $s \in \mathbb{R} \backslash\{0\}$; the matrix $Q_{1}^{s}$ is degenerate if and only if $-\frac{1}{s} \in \operatorname{Sp}(K)$; moreover, dim $\operatorname{ker} Q_{1}^{s}$ equals the multiplicity of the eigenvalue $-\frac{1}{s}$.

Proof. Let $v \in \mathcal{V}$; it is easy to see that

$$
\left\langle J Z_{\tau}^{s} v(\tau), Z_{t}^{s} v(t)\right\rangle=s\left\langle J Z_{\tau} v(\tau), Z_{t} v(t)\right\rangle
$$

Hence

$$
<(I+s K) v \mid v>=\int_{0}^{1}\langle v(t), v(t)\rangle d t-\int_{0}^{1}\left\langle J \int_{0}^{t} Z_{\tau}^{s} v(\tau) d \tau, Z_{t}^{s} v(t)\right\rangle d t
$$

We keep symbol $\mathcal{V}$ for the complexification of $\mathcal{V}$, i. e. $\mathcal{V}=\mathbb{C} \otimes \mathcal{V}$. A vectorfunction $v \in \mathcal{V}$ is an eigenvector of $K$ with the eigenvalue $-\frac{1}{s}$ if and only if

$$
\begin{equation*}
\int_{0}^{1}\left\langle v(t)-Z_{t}^{s *} J \int_{0}^{t} Z_{t}^{s} v(\tau) d \tau, w(t)\right\rangle d t=0, \quad \forall w \in \mathcal{V} \tag{17}
\end{equation*}
$$

where $\mathcal{V}=\left\{w \in L_{2}\left([0,1] ; \mathbb{C}^{d}\right): \int_{0}^{1} X_{t} w(t) d t=0\right\}$.
Equation $\int_{0}^{1} X_{t} w(t) d t=0$ can be rewritten as follows: $\int_{0}^{1}\left\langle J \nu, Z_{t} w(t)\right\rangle d t=$ 0 for any $\nu \in \operatorname{ker} \pi_{*}$. In other words, $\mathcal{V}^{\perp}=\left\{t \mapsto Z_{t}^{s *} J \nu: \nu \in \operatorname{ker} \pi_{*}\right\}$. Hence relation (17) is equivalent to the existence of $\nu \in \operatorname{ker} \pi_{*}$ such that

$$
v(t)=Z_{t}^{s *} J\left(\int_{0}^{t} Z_{\tau}^{s} v(\tau) d \tau+\nu\right)
$$

Moreover, the vector $\nu$ is unique for given vector-function $v$ since the Volterra equation $v(t)=Z_{t}^{s *} J \int_{0}^{t} Z_{\tau}^{s} v(\tau) d \tau$ has only zero solution. We set $\eta(t)=$ $\int_{0}^{t} Z_{\tau}^{s} v(\tau) d \tau+\nu$ and obtain that a vector-function $v$ satisfies relation (17) if and only if

$$
\begin{equation*}
\dot{\eta}=Z_{t}^{s} Z_{t}^{s *} J \eta, \quad 0 \leq t \leq 1, \quad \eta(0), \eta(1) \in \operatorname{ker} \pi_{*} . \tag{18}
\end{equation*}
$$

It follows that $\operatorname{dim} \operatorname{ker} Q_{1}^{s}$ is equal to the multiplicity of the eigenvalue $-\frac{1}{s}$ of the operator $K$ plus the dimension of the space of constant solutions of equation (18) that belong to ker $\pi_{*}$. Let $(\nu, 0) \in \mathbb{C}^{d} \times \mathbb{C}^{d}$ be such a solution. We plug-in it in (18) and get $X_{t} X_{t}^{*} \nu=0, \forall t \in[0,1]$. Hence $\nu=0$.

Corollary 2. The equation $\operatorname{det} Q_{1}^{s}=0, s \in \mathbb{C}$, has only real roots.
Indeed, the operator $K$ is symmetric and has only real eigenvalues.
Proposition 6. Let $s_{0} \in \mathbb{R} \backslash\{0\}$ and $\operatorname{det} Q_{1}^{s_{0}}=0$. Then the multiplicity of the root $s_{0}$ of the equation $\operatorname{det} Q_{1}^{s}=0$ is equal to $\operatorname{dim} \operatorname{ker} Q_{1}^{s_{0}}$.

Proof. We may assume that $s \in \mathbb{R}$ and work in the real setting. We denote by $\eta^{s}(t)=\left(p^{s}(t), q^{s}(t)\right), t \in[0,1]$, those solutions of the Hamiltonian system

$$
\begin{equation*}
\dot{\eta}^{s}=s Z_{t} Z_{t}^{*} J \eta^{s} \tag{17}
\end{equation*}
$$

that satisfy the initial condition $q^{s}(0)=0$; then $q^{s}(t)=s Q_{t}^{s} p^{s}(0)$ (c. f. proof of Proposition 4).

Let $\Lambda_{1}^{s}=\left\{\left(p^{s}(1), q^{s}(1)\right): p^{s}(0) \in \mathbb{R}^{d}\right\}$; then $\Lambda_{1}^{s}$ is a Lagrange subspace of the symplectic space $\mathbb{R}^{d} \times \mathbb{R}^{d}$ endowed with the standard symplectic structure $\left(\left(p_{1}, q_{1}\right),\left(p_{2}, q_{2}\right)\right) \mapsto\left\langle p_{1}, q_{2}\right\rangle-\left\langle p_{2}, q_{1}\right\rangle$.

We take another Lagrange subspace $\Delta$ that is transversal to $\operatorname{ker} \pi_{0}^{*}$ and $\Lambda_{1}^{*}$. Then $\Delta=\left\{(A q, q): q \in \mathbb{R}^{d}\right\}$ where $A$ is a symmetric matrix. We make a symplectic change of variables $(p, q) \mapsto\left(p^{\prime}, q\right)$ by putting $p^{\prime}=p-A q$. In new variables, $\Lambda_{1}^{s}$ is transversal to the "horizontal" subspace defined by the equation $p^{\prime}=0$ for all $s$ close to $s_{0}$. Hence

$$
\Lambda_{1}^{s}=\left\{\left(p^{\prime}, R(s) p^{\prime}\right): p^{\prime} \in \mathbb{R}^{d}\right\}
$$

where $R(s)$ is a symmetric matrix. According to our construction, $s Q_{1}^{s}$ is the product of an analytic with respect to $s$ nondegenerate matrix and $R(s)$. It follows that $\operatorname{dim} \operatorname{ker} R(s)=\operatorname{dim} \operatorname{ker} Q_{1}^{s}$ and the multiplicities of the root $s_{0}$ of the equations $\operatorname{det} R(s)=0$ and $\operatorname{det} Q_{1}^{s}=0$ are equal. We denote this multiplicity by $\mu\left(s_{0}\right)$.

Standard perturbation theory for the eigenvalues of symmetric operators (see, for instance, [5]) implies that: (i) $\mu\left(s_{0}\right) \geq \operatorname{dim} \operatorname{ker} R\left(s_{0}\right)$; (ii) if $\operatorname{det} \frac{d R}{d s}\left(s_{0}\right) \neq 0$, then $\mu\left(s_{0}\right)=\operatorname{dim} \operatorname{ker} R\left(s_{0}\right)$.

Now we have to compute $\frac{d R}{d s}$. To do that, we re-write equation (17) in coordinates $\left(p^{\prime}, q\right)$. In new coordinates, this Hamiltonian equation has the same structure as (17) with $Z_{t}=\left(Y_{t}, X_{t}\right)$ substituted by $Z_{t}^{\prime}=\left(Y_{t}+A X_{t}, X_{t}\right)$. From now on, we'll write $p, Z_{t}$ instead of $p^{\prime}, Z_{t}^{\prime}$ and work directly with equation (17) in order to simplify notations. We have $\eta^{s}(1)=\left(p^{s}(1), R(s) q^{s}(1)\right)$; the symmetricity of the matrix $R(s)$ implies that

$$
\left\langle J \eta^{s}(1), \frac{\partial \eta^{s}}{\partial s}(1)\right\rangle=\left\langle p^{s}(1),\left(\frac{d R}{d s}\right) p^{s}(1)\right\rangle .
$$

Proposition 6 is now reduced to the following lemma:
Lemma 1. If $\left\langle J \eta^{s}(1), \frac{\partial \eta^{s}}{\partial s}(1)\right\rangle=0$, then $\eta^{s}(1)=0$.
Proof. We have: $\frac{\partial \dot{\eta}^{s}}{\partial s}=s Z_{t} Z_{t}^{*} J \frac{\partial \eta^{s}}{\partial s}+Z_{t} Z_{t}^{*} J \eta^{s}$. Hence

$$
\begin{aligned}
\frac{\partial}{\partial t}\left\langle J \eta^{s}(t), \frac{\partial \eta^{s}}{\partial s}(t)\right\rangle=s\langle & \left.J Z_{t} Z_{t}^{*} J \eta^{s}(t), \frac{\partial \eta^{s}}{\partial s}(t)\right\rangle+s\left\langle J \eta^{s}(t), Z_{t} Z_{t}^{*} J \frac{\partial \eta^{s}}{\partial s}(t)\right\rangle \\
& +\left\langle J \eta^{s}(t), Z_{t} Z_{t}^{*} J \eta^{s}\right\rangle .
\end{aligned}
$$

First two terms in the right hand side of the last identity cancel because the matrix $J$ is anti-symmetric and the matrix $J Z_{t} Z_{t}^{*} J$ is symmetric. We obtain:

$$
\frac{\partial}{\partial t}\left\langle J \eta^{s}(t), \frac{\partial \eta^{s}}{\partial s}(t)\right\rangle=\left\langle Z_{t}^{*} J \eta^{s}(t), Z_{t}^{*} J \eta^{s}(t)\right\rangle
$$

Moreover, $\left\langle J \eta^{s}(0), \frac{\partial \eta^{s}}{\partial s}(0)\right\rangle=0$ since $q^{s}(0)=0$ for all $s$. We get:

$$
\left\langle J \eta^{s}(1), \frac{\partial \eta^{s}}{\partial s}(1)\right\rangle=\int_{0}^{1}\left|Z_{t}^{*} J \eta^{s}(t)\right|^{2} d t
$$

If $\left\langle J \eta^{s}(1), \frac{\partial \eta^{s}}{\partial s}(1)\right\rangle=0$, then $Z_{t}^{*} J \eta^{s}(t)=0, \forall t \in[0,1]$. Equation (17) now implies that $\dot{\eta}^{s}(t) \equiv 0$ and $\eta^{s}(t)=\left(p^{s}(0), 0\right)$. We plug-in this equilibrium in (17) and obtain: $X_{t} X_{t}^{*} p^{s}(0)=0,0 \leq t \leq 1$. Hence $\Gamma_{1} p^{s}(0)=$ $\int_{0}^{1} X_{t} X_{t}^{*} d t p^{s}(0)=0$. Recall that $\Gamma_{1}$ is a nondegenerate matrix.

Proposition 4 implies that $s \mapsto \operatorname{det} Q_{1}^{s}, s \in \mathbb{C}$, is an entire function that satisfies the estimate $\operatorname{det} Q_{1}^{s} \leq c e^{c|s|}$ for some constant $c$. It follows from a classical Hadamard theorem (see [3]) that such a function has a presentation $\operatorname{det} Q_{1}^{s}=a e^{b s} \prod_{i}\left(1-\frac{s}{s_{i}}\right)$, where $s_{i}$ are roots of the equation $\operatorname{det} Q_{1}^{s}=0$ counted according to multiplicity and $a, b \in \mathbb{C}$.

Propositions 5 and 6 now imply that

$$
\begin{equation*}
\operatorname{det} Q_{1}^{s}=a e^{b s} \operatorname{det}(I+s K), \quad s \in \mathbb{C} . \tag{18}
\end{equation*}
$$

It remains to find $a$ and $b$. We have: $\operatorname{det} Q_{1}^{0}=a$. Moreover,

$$
\left.\frac{d}{d s} \operatorname{det} Q_{1}^{s}\right|_{s=0}=\operatorname{tr}\left(\left.\frac{d Q_{1}^{s}}{d s}\right|_{s=0}\left(Q_{1}^{0}\right)^{-1}\right) \operatorname{det} Q_{1}^{0}
$$

We differentiate identity (18) and obtain: $\operatorname{tr}\left(\left.\frac{d Q_{1}^{s}}{d s}\right|_{s=0}\left(Q_{1}^{0}\right)^{-1}\right)=b+\operatorname{tr} K$.
Let $\Pi: L^{2}\left([0,1] ; \mathbb{R}^{m}\right) \rightarrow \mathcal{V}$ be the orthogonal projection, then $\operatorname{tr} K=$ $\operatorname{tr}(\Pi \hat{K} \Pi)=\operatorname{tr}(K \Pi)$, where $\hat{K} v(t)=-\int_{0}^{t} Z_{t}^{*} J Z_{\tau} v(\tau) d \tau$. The subspace $\mathcal{V}$ is the kernel of the operator $A: v \mapsto \int_{0}^{1} X_{t} v(t) d t, v \in \mathrm{~L}^{2}\left([0,1] ; \mathbb{R}^{m}\right)$. We have: $\quad \Pi=I-A^{*}\left(A A^{*}\right)^{-1} A, \quad \Pi v(t)=v(t)-X_{t}^{*} \Gamma_{1}^{-1} \int_{0}^{1} X_{\tau} v(\tau) d \tau$.

The trace of a trace-class operator $\bar{B}$ of a form $(\bar{B} v)(t)=\int_{0}^{1} B(t, \tau) v(\tau) d \tau$ is computed according to the formula $\operatorname{tr} \bar{B}=\int_{0}^{1} \operatorname{tr} B(t, t) d t$. This presentation of the trace is valid in our situation as well. Moreover, $\operatorname{tr}\left(Z_{t}^{*} J Z_{t}\right) \equiv 0$; we collect the terms and obtain the formula:

$$
\begin{aligned}
& \operatorname{tr} K=\int_{0}^{1} \int_{0}^{t} \operatorname{tr}\left(Z_{t}^{*} J Z_{\tau} X_{\tau}^{*} \Gamma_{1}^{-1} X_{t}\right) d \tau d t \\
& \quad=\int_{0}^{1} \int_{0}^{t} \operatorname{tr}\left(X_{t} Z_{t}^{*} J Z_{\tau} X_{\tau}^{*} \Gamma_{1}^{-1}\right) d \tau d t
\end{aligned}
$$

as in the statement of Theorem 2 (recall that we are working in coordinated where $\left.H_{t}=-I\right)$.

Now we compute $Q_{1}^{0}$. For $s=0$, system (18) is reduced to the system: $\dot{p}=0, \dot{q}=X_{t} X_{t}^{*} p$. Hence $Q_{1}^{0}=\int_{0}^{1} X_{t} X_{t}^{*} d t=\Gamma_{1}$.

The last step is the calculation of $\left.\frac{d Q_{1}^{s}}{d s}\right|_{s=0}$. System (16) reads:

$$
\begin{gathered}
\dot{p}=s Y_{t} X_{t}^{*} p-s^{2} Y_{t} Y_{t}^{*} q \\
\dot{q}=X_{t} X_{t}^{*} p-s X_{t} Y_{t}^{*} q .
\end{gathered}
$$

We have:

$$
q^{s}(t)=\int_{0}^{t} X_{\tau} X_{\tau}^{*} d \tau p_{0}+O(s), \quad p^{s}(t)=s \int_{0}^{t} Y_{\tau} X_{\tau}^{*} d \tau p_{0}+O\left(s^{2}\right) .
$$

Hence

$$
\begin{gathered}
q^{s}(1)=\int_{0}^{1} X_{t} X_{t}^{*} p(t)^{s}-s X_{t} Y_{t}^{*} q^{s}(t) d t= \\
\Gamma_{1} p_{0}+s \int_{0}^{1} \int_{0}^{t} X_{t} X_{t}^{*} Y_{\tau} X_{\tau}^{*}-X_{t} Y_{t}^{*} X_{\tau} X_{\tau}^{*} d \tau d t p_{0}+O\left(s^{2}\right)
\end{gathered}
$$

Moreover, $X_{t} X_{t}^{*} Y_{\tau} X_{\tau}^{*}-X_{t} Y_{t}^{*} X_{\tau} X_{\tau}^{*}=X_{t} Z_{t}^{*} J Z_{\tau} X_{\tau}^{*}$. We see that $\operatorname{tr}\left(\left.\frac{d Q_{1}^{s}}{d s}\right|_{s=0} \Gamma_{1}^{-1}\right)=\operatorname{tr} K$. Hence $b=0$.

Remark 5. We used only the first derivative of $Q_{1}^{s}$ at $s=0$ in the proof of Theorem 2 but system (16) allows us to find explicit integral expressions for all higher derivatives and thus to obtain integral expressions for all elementary symmetric functions of the eigenvalues of the operator $K$.

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[^1]:    ${ }^{1}$ We say that a self-adjoint map is positive (negative) if the corresponding quadratic form is positive (negative).

