# Non Abelian Asymptotics of Szegö Kernels on Compact Quantized Manifolds 

Doctoral Thesis



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#### Abstract

Let $M$ be a complex projective manifold with a positive line bundle $A$ on it. The unit circle bundle $X$, inside the dual of $A$, is a contact and CR manifold by positivity of $A$. Let $H(X) \subset L^{2}(X)$ be the associated Hardy space, and $\Pi$ be the corresponding orthogonal projector, which is known in the literature as the Szegö projector. Assume given a Hamiltonian and holomorphic action of a compact Lie group $G$ on $M$, and suppose that the action lifts to a contact CR action on $X$. There is a naturally induced unitary representation of $G$ on $H(X)$, which therefore splits equivariantly into a direct sum of isotypical components, indexed by the irreducible representations of G. Assuming that the moment map is nowhere vanishing, each isotypical summand is finite-dimensional, and therefore the kernel of the corresponding orthogonal projector ('equivariant Szegö kernel') is a $C^{\infty}$ function on $X \times X$. One is led to study the local asymptotics of the equivariant projectors pertaining to the irreducible representations in a given ray in weight space. In this thesis we consider the case of special unitary group $S U(2)$ and unitary group $U(2)$.


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## Contents

1 Introduction ..... 4
1.1 Geometric setting and goal ..... 4
1.2 The case $G=S U(2)$ ..... 7
1.3 The case $G=U(2)$ ..... 11
1.4 Recalls on prior literature ..... 18
2 Examples ..... 23
2.1 Examples about $G=S U(2)$ ..... 23
2.2 Example about $G=U(2)$ ..... 26
3 Preliminaries ..... 34
3.1 Quantized manifolds ..... 34
3.2 Compact Lie groups and actions ..... 35
3.2.1 Volume elements ..... 35
3.2.2 Harmonic analysis on compact Lie groups ..... 35
3.2.3 Hamiltonian group actions ..... 38
3.3 The Szegö projector ..... 39
3.3.1 The Szegö kernel ..... 39
3.3.2 Equivariant Szegö projector ..... 40
4 Proofs ..... 42
4.1 Proof of Theorem 1.2.1 ..... 42
4.2 Proof of Theorem 1.2.2 ..... 49
4.3 Proof of Theorem 1.2.3 ..... 56
4.4 Proof of Theorem 1.3.1 ..... 67
4.4.1 The proof ..... 70
4.5 Proof of Theorem 1.3.2 ..... 71
4.5.1 Preliminaries ..... 71
4.5.2 The proof ..... 74
4.6 Proof of Theorem 1.3.3 ..... 80
4.6.1 An a priori polynomial bound ..... 80
4.6.2 The proof ..... 81
4.7 Proof of Theorem 1.3.4, 1.3.6 and 1.3.7 ..... 91
4.7.1 Preliminaries on local rescaled asymptotics ..... 91
4.7.2 Proof of Theorem 1.3.4 ..... 99
4.7.3 Proof of Theorem 1.3.6 ..... 100
4.7.4 Proof of Theorem 1.3.7 ..... 102
4.8 Proof of Theorem 1.3.5 ..... 104

## Chapter 1

## Introduction

### 1.1 Geometric setting and goal

Let $M$ be a connected compact complex $d$-dimensional projective manifold, with a complex structure $J$ (for a more detailed discussion and for definitions see the preliminaries in the next chapter). Let $(A, h)$ be a positive line bundle on $M$, where $h$ is the Hermitian structure, with connection $\nabla$ compatible with the metric such that curv $\nabla=-2 i \omega$, where $\omega$ is a Kähler form on $M$. We shall denote by $g$ the corresponding Riemannian structure on $M$, given by

$$
\begin{equation*}
g_{m}(\mathbf{v}, \mathbf{w}):=\omega_{m}\left(\mathbf{v}, J_{m}(\mathbf{w})\right) \quad\left(m \in M, \mathbf{v}, \mathbf{w} \in T_{m} M\right) \tag{1.1}
\end{equation*}
$$

Furthermore $M$ has a natural choice of a volume form: $\mathrm{d} V_{M}:=(1 / d!) \omega^{d}$.
If $A^{\vee} \supset X \xrightarrow{\pi} M$ is the unit circle bundle in the dual of $A$, then $\nabla$ naturally corresponds to a connection 1-form $\alpha$ on $X$, such that $\mathrm{d} \alpha=2 \pi^{*}(\omega)$. Let us notice that $(X, \alpha)$ is a contact manifold and inherits the volume form $\mathrm{d} \mathrm{V}_{X}:=(2 \pi)^{-1} \alpha \wedge \pi^{*}\left(\mathrm{~d} \mathrm{~V}_{M}\right)$. Furthermore $\alpha$ determines an invariant splitting of the tangent bundle of $X$ as

$$
\begin{equation*}
T X:=\mathcal{V}(X / M) \oplus \mathcal{H}(X / M) \tag{1.2}
\end{equation*}
$$

where $\mathcal{V}(X / M):=\operatorname{ker}(\mathrm{d} \pi)$ is the vertical tangent bundle, and $\mathcal{H}(X / M):=$ $\operatorname{ker}(\alpha)$ is the horizontal one. Given $V \in \mathfrak{X}(M)$ (the Lie algebra of smooth vector fields on $M$ ), we shall denote by $V^{\sharp} \in \mathfrak{X}(X)$ its horizontal lift to $X$. If the vector field $\partial_{\theta} \in \mathfrak{X}(X)$ is the generator of the structure $S^{1}$-action, then $\partial_{\theta}$ spans the vertical tangent bundle.

The holomorphic structure on $M$, pulled-back to the horizontal tangent bundle, endows $X$ with a CR structure. Explicitly, the complex structure $J$ on $M$ naturally lifts to a vector bundle endomorphism of $T X$, also denoted by $J$, such that $J\left(\partial_{\theta}\right)=0$ and

$$
\begin{equation*}
J\left(\mathbf{v}^{\sharp}\right)=J(\mathbf{v})^{\sharp} \quad(\mathbf{v} \in \mathfrak{X}(M)) . \tag{1.3}
\end{equation*}
$$

By the positivity of $(A, h), X$ is the boundary of a strictly pseudo-convex domain. The corresponding Hardy space $H(X) \subset L^{2}(X)$ encapsulates the holomorphic structure of $A$ and its tensor powers. The corresponding orthogonal projector $\Pi: L^{2}(X) \rightarrow H(X)$ is called the Szegö projector. For more details concerning the Szegö Kernel and its description as a Fourier Integral Operator see section 3.3.1 in the next chapter.

Suppose that, in addition, one is given a compact Lie group $G$ and a holomorphic Hamiltonian action $\mu: G \times M \rightarrow M$ with moment map $\Phi_{G}$ : $M \rightarrow \mathfrak{g}^{\vee}$; then $G$ acts on $M$ via Riemannian isometries (see the next chapter for a more detailed description of the representation theory of compact Lie groups and for Hamiltonian group actions). In particular, for each $\xi \in \mathfrak{g}$, let us denote with $\xi_{M} \in \mathfrak{X}(M)$ the vector field associated to $\xi$. We shall assume that the action $\mu$ lifts to an action $\tilde{\mu}: G \times X \rightarrow X$ preserving the contact and CR structures. The vector field $\xi_{X}$ associated to $\xi \in \mathfrak{g}$ is related to $\xi_{M}$ via the following formula

$$
\begin{equation*}
\xi_{X}(x)=\xi_{M}^{\sharp}(m)-\left\langle\Phi_{G}(m), \xi\right\rangle \partial_{\theta} . \tag{1.4}
\end{equation*}
$$

Then pull-back of functions, given by $g \cdot s:=\tilde{\mu}_{g^{-1}}^{*}(s)$, is a unitary representation of $G$ on $L^{2}(X)$ leaving $H(X) \subset L^{2}(X)$ invariant. This yields a unitary representation $\hat{\mu}: G \rightarrow U(H(X))$ which commutes with the standard $S^{1}$ action.

By the Theorem and Peter and Weyl, all the irreducible representations of $G$ are finite dimensional, and any unitary representation of $G$ splits equivariantly as a direct sum of irreducible representations of $G$.

We shall focus in particular on the two cases $G=S U(2)$ and $G=U(2)$. If $G=U(2)$, the irreducible representations of $G$ are indexed by the couples $\boldsymbol{\nu}=\left(\nu_{1}, \nu_{2}\right)$ of integers satisfying $\nu_{1}>\nu_{2}$. There is an equivariant unitary isomorphism

$$
H(X) \cong \bigoplus_{\nu_{1}>\nu_{2}} H(X)_{\nu}
$$

where $H(X)_{\boldsymbol{\nu}} \subset H(X)$ is the $\boldsymbol{\nu}$-isotypical component. Correspondingly,

$$
\Pi=\sum_{\nu_{1}>\nu_{2}} \Pi_{\nu}
$$

where $\Pi_{\nu}: L^{2}(X) \rightarrow H(X)_{\nu}$ is the orthogonal projector. In the same spirit, the action of $S U(2)$ gives rise to the decomposition

$$
H(X) \cong \bigoplus_{\nu>0} H(X)_{\nu}
$$

where here $\boldsymbol{\nu}=(\nu, 0)$ and the corresponding equivariant Szegö kernels will be denoted by $\Pi_{\nu}$.

In general, $H(X)_{\nu}$ may well be infinite dimensional; however, if $\mathbf{0} \notin$ $\Phi_{G}(M)$ then $\operatorname{dim}\left(H(X)_{\nu}\right)<+\infty$, see [Pao12]. In this case $\Pi_{\nu}$ is a smoothing operator, with a distributional kernel

$$
\Pi_{\nu}(\cdot, \cdot) \in \mathcal{C}^{\infty}(X \times X)
$$

In particular,

$$
\operatorname{dim} H(X)_{\nu}=\int_{X} \Pi_{\nu}(x, x) \mathrm{d} V_{X}(x)
$$

The goal of this thesis is to study the concentration behavior of $\Pi_{k \nu}$ when $k \rightarrow+\infty$, in the specific cases where $G=S U(2)$ and $G=U(2)$.

A few words about the organization of this thesis. In chapter $\S 3$ we review some definitions and theorems about harmonic analysis, moment maps and Szegö projector.

In chapter $\S 1$ we described the main results of this thesis. In section 1.2 (based on [GP18a]) we study the case $G=S U(2)$. We obtain an asymptotic expansion in descending powers of $k$ for both the free and locally free action on $X$. First we provide some quantitative information about the rate of decay of $\Pi_{k \nu}(x, y)$ when $x$ approaches to $G \cdot y$. An explicit expression of the leading order term is found for the asymptotic of $\Pi_{k \nu}(x, x)$ for a locally free action at $x$ and for a near-diagonal rescaled displacement when the the stabilizer of $x$ is contained in the center of $G$. In section 1.3 (based on [GP18b]) we prove the results for $U(2)$. It turns out that a very special role is played by a locus $M_{\mathcal{O}_{\nu}}^{G}$, which is the inverse image via the moment map $\Phi_{G}$ of the cone through the co-adjoint orbit of $\boldsymbol{\nu}$. Under some reasonable hypotheses,
$M_{\mathcal{O}_{\nu}}^{G}$ is an hyper-surface and it divides $M$ in two connected components: the "outside" and the "inside". We are able to give a precise result about the rate of decay of the asymptotics of $\Pi_{k \nu}(x, y)$ when one approaches to $M_{\mathcal{O}_{\nu}}^{G}$ from the outside. Nevertheless, we prove that the equivariant Szegö kernel has a rapidly decreasing asymptotic on compact subset not lying in $\mathcal{Z}_{\nu}:=$ $\left\{(x, y) \in X_{\mathcal{O}_{\nu}}^{G} \times X_{\mathcal{O}_{\nu}}^{G}: y \in G \cdot x\right\}$, where $X_{\mathcal{O}_{\nu}}^{G}$ denotes the pull-back of $M_{\mathcal{O}_{\nu}}^{G}$ on $X$. If the action of $G$ on $X_{\mathcal{O}_{\nu}}^{G}$ is free, we find an explicit expression for the asymptotic along the diagonal, an asymptotic expansion for near-diagonal rescaled displacements and a lower bound for the dimension of the isotypes $H_{k \nu}(X)$. Similar results are given also for the locally free case.

The proofs are collected in the last chapter. The case of $G=S U(2)$ bears a close resemblance to the case of $S^{1}$, in sharp contrast with the case of $G=U(2)$; this reflects the fact that the cone over the coadjoint orbit is open and dense in the case of $S U(2)$ and $S^{1}$, but has codimension 1 in the case of $U(2)$. The first main step in the proofs is to show that the equivariant Szegö kernel has the same asymptotic as a compact supported oscillatory integral. The leading term of the asymptotic expansion is found by the use of stationary phase Lemma. In the case $G=U(2)$, the proof concerning the rapidly decreasing asymptotic of $\Pi_{k \nu}(x, y)$ on compact subset in the complement of $\mathcal{Z}_{\boldsymbol{\nu}}$ has a different approach: it is grounded on results of Guillemin and Sternberg ([GS82c]) regarding functorial properties of some distributions related with the equivariant Szegö kernel.

### 1.2 The case $G=S U(2)$

In this section we denote by $G$ the group $S U(2)$, The Lie algebra $\mathfrak{g}$ of $S U(2)$ consists of skew-Hermitian matrices with trace zero. There is a natural $G$ invariant Euclidean scalar product on $\mathfrak{g}$, given by $\left\langle\xi_{1}, \xi_{2}\right\rangle=\operatorname{trace}\left(\xi_{1} \bar{\xi}_{2}^{t}\right)$. By means of $\langle\cdot, \cdot\rangle$, we can equivariantly identify $\mathfrak{g}$ and $\mathfrak{g}^{\vee}$. The first step will be to provide quantitative information about the rate of decay of $\Pi_{k \nu}(x, y)$ when $x$ approaches to $G \cdot y$.

Theorem 1.2.1. Assume that $\mathbf{0}$ does not lie in the image of the moment map $\Phi_{G}(M)$. Let $C, \epsilon>0$. Then, uniformly for

$$
\operatorname{dist}_{X}(x, G \cdot y) \geq C k^{\epsilon-1 / 2}
$$

we have $\Pi_{k \nu}(x, y)=O\left(k^{-\infty}\right)$ as $k \rightarrow+\infty$.

Before stating our next Theorem, we need to introduce some terminology. Let us suppose that $\mathbf{0} \notin \Phi_{G}(M)$. Thus there exists one and only one $h_{m} T \in$ $G / T$ such that

$$
h_{m} \Phi_{G}(m) h_{m}^{-1}=i\left(\begin{array}{cc}
\lambda(m) & 0 \\
0 & -\lambda(m)
\end{array}\right)
$$

where $\lambda(m)$ is the unique positive eigenvalue of $\Phi_{G}(m)$ and clearly $\lambda: M \rightarrow$ $(0,+\infty)$ is smooth.

Remark 1. The positive eigenvalue $\lambda(m)$ has a symplectic interpretation, being closely related to the moment map for the action restricted to a suitable torus $T_{m} \leq G$. Let us set

$$
\beta:=\left(\begin{array}{cc}
i & 0 \\
0 & -i
\end{array}\right) ;
$$

thus $\beta$ is the infinitesimal generator of the standard torus $T$, and therefore $\operatorname{Ad}_{h_{m}}(\beta)$ is the infinitesimal generator of the torus $T_{m}:=C_{h_{m}}(T)$ (here $C_{g}(h):=g h g^{-1}$, for all $\left.g, h \in G\right)$. Then for any $m \in M$ we have

$$
2 \lambda(m)=\left\langle h_{m}^{-1} \Phi_{G}(m) h_{m}, \beta\right\rangle=\left\langle\Phi_{G}(m), \operatorname{Ad}_{h_{m}}(\beta)\right\rangle
$$

Let us denote by $G_{x}$ the stabilizer of $x \in X$. By equivariance of $\Phi_{G}, G_{x}$ stabilizes $\Phi_{G}(m)$. Hence, $G_{x} \subset h_{m_{x}} T h_{m_{x}}^{-1}$. In particular, if the lifted action $\tilde{\mu}: G \times X \rightarrow X$ is locally free at $x$, then $G_{x}$ is finite and Abelian. There exists $e^{i \vartheta_{j}} \in S^{1}, j=1, \ldots, N_{x}$, such that

$$
G_{x}=\left\{h_{m_{x}}\left(\begin{array}{cc}
e^{i \vartheta_{j}} & 0 \\
0 & e^{-i \vartheta_{j}}
\end{array}\right) h_{m_{x}}^{-1}: \quad j=1, \ldots, N_{x}\right\} .
$$

Let $Z:=\left\{ \pm I_{2}\right\}$ be the center of $G$, and set $Z_{x}:=G_{x} \cap Z$. We shall see that each $g \in G_{x}$ contributes to the rescaled asymptotics of $\Pi_{k \nu}$ near $x$, and that the nature of the contribution is quite different depending on whether $g \in Z_{x}$ or $g \in G \backslash G_{x}$.

Definition 1.2.1. If $l \in \mathbb{Z}$, let us define $f_{l}: T \rightarrow \mathbb{C}$ by letting

$$
f_{l}: e^{i \vartheta} I_{2} \in T \mapsto e^{i l \vartheta} \in \mathbb{C}^{*} .
$$

If $h \in G_{x} \backslash Z_{x}$, then $h \neq h^{-1}$. Hence $G_{x} \backslash Z_{x}$ has even cardinality $b_{x}=2 a_{x}$, and perhaps after renumbering its elements can be arranged in pairs $\left(g_{j}, g_{j}^{-1}\right)$, $j=1, \ldots a_{x}$. For every $j=1, \ldots, a_{x}$. For every $j=1, \ldots, a_{x}$, let us set

$$
t_{j}:=h_{m_{x}}^{-1} g_{j} h_{m_{x}}=\left(\begin{array}{cc}
e^{i \vartheta_{j}} & 0  \tag{1.5}\\
0 & e^{-i \vartheta_{j}}
\end{array}\right) .
$$

Definition 1.2.2. If $z \in \mathbb{C}$, let us set

$$
A(z):=i\left(\begin{array}{cc}
0 & z \\
\bar{z} & 0
\end{array}\right) \in \mathfrak{g} .
$$

The the $\mathbb{R}$-linear map

$$
\eta_{j}: z \in \mathbb{C} \mapsto\left(\operatorname{Ad}_{t_{j}^{-1}}-\operatorname{id}_{\mathfrak{g}}\right)(A(z)) \in \mathfrak{g}
$$

is injective. Therefore, since $\tilde{\mu}$ is locally free at $x$, there is a positive $2 \times 2$ matrix $C(x ; j)$ such that

$$
\left\|\operatorname{Ad}_{h_{m_{x}}}\left(\eta_{j}(z)\right)_{X}(x)\right\|^{2}=\frac{1}{2} \cdot Z^{t} C(x ; j) Z \quad(z \in \mathbb{C})
$$

where $Z:=(a, b)^{t} \in \mathbb{R}^{2}$ if $z=a+i b$. Let us define

$$
B(x ; j):=C(x ; j)+4 i \sin \left(2 \vartheta_{j}\right) \cdot \lambda\left(m_{x}\right) I_{2} .
$$

Finally, let us denote by $V_{3}$ the area of the unit sphere $S^{3} \subset \mathbb{R}^{4}$, and set

$$
D_{G / T}:=2 \pi / V_{3}
$$

We shall prove the following.
Theorem 1.2.2. Le us assume that $\mathbf{0} \notin \Phi_{G}(M), \tilde{\mu}$ is locally free at $x$ and that $G_{x} \backslash Z_{x}=\left\{g_{1}, g_{1}^{-1}, \ldots, g_{a_{x}}, g_{a_{x}}^{-1}\right\}$. Then as $k \rightarrow+\infty$ there is an asymptotic expansion

$$
\Pi_{k \nu}(x, x) \sim \Pi_{k \nu}(x, x)_{Z_{x}}+\Pi_{k \nu}(x, x)_{G_{x} \backslash Z_{x}}
$$

where

$$
\begin{aligned}
\Pi_{k \nu}(x, x)_{Z_{x}} \sim & \frac{1}{2 \lambda\left(m_{x}\right)} \cdot\left(\frac{\nu k}{2 \pi \lambda\left(m_{x}\right)}\right)^{d} \cdot \sum_{g \in Z_{x}} f_{1-k \cdot \nu}(g) \\
& \cdot\left[1+\sum_{j=1}^{+\infty} k^{-j / 2} B_{g j}(x)\right]
\end{aligned}
$$

and

$$
\begin{aligned}
\Pi_{k \nu}(x, x)_{G_{x} \backslash Z_{x}} \sim & 4 \pi \cdot D_{G / T} \cdot\left(\frac{\nu k}{2 \pi \lambda\left(m_{x}\right)}\right)^{d} \\
& \cdot\left[\sum_{j=1}^{a_{x}} \Re\left(\frac{i \sin \left(\vartheta_{j}\right) \cdot e^{-i k \nu \cdot \vartheta_{j}}}{\sqrt{\operatorname{det}(B(x ; j))}}\right)+\sum_{l \geq 1} k^{-l / 2} P_{j l}\left(m_{x}\right)\right]
\end{aligned}
$$

for appropriate smooth functions $B_{g j}, P_{j l}: M \rightarrow \mathbb{R}$.

We shall consider the asymptotic expansion for near-diagonal rescaled displacement. Let us analyze the less general case where $\tilde{\mu}$ is locally free at $x$ and $G_{x}$ is contained in $Z$.

Let us choose a set of Heisenberg local coordinates centered at $x \in X$ for which we refer to [SZ02]. We shall set $x+(\theta, \mathbf{v})$ for $v=(\theta, \mathbf{v}) \in(-\pi, \pi) \times \mathbb{R}^{2 d}$. When $\theta=0$, we shall write $x+\mathbf{v}$ for $x+(0, \mathbf{v})$.

Definition 1.2.3. If $m \in M$ and $\mathbf{v}_{1}, \mathbf{v}_{2} \in T_{m} M$, following [SZ02] let us set

$$
\psi_{2}\left(\mathbf{v}_{1}, \mathbf{v}_{2}\right):=-i \omega_{m}\left(\mathbf{v}_{1}, \mathbf{v}_{2}\right)-\frac{1}{2}\left\|\mathbf{v}_{1}-\mathbf{v}_{2}\right\|^{2}
$$

Let us set $\mathbf{v}^{(g)}:=\mathrm{d}_{m_{x}} \mu_{g}(\mathbf{v})$ and define

$$
u_{0}(\nu, x):=\frac{\nu}{2 \lambda\left(m_{x}\right)} .
$$

Given this definition we can formulate our result about the scaling asymptotics for $\Pi_{k \nu}(x, x)$.

Theorem 1.2.3. Let us assume that $\mathbf{0} \notin \Phi_{G}(M)$ and that $\tilde{\mu}$ is locally free on $X$. Let $G_{x}$ be the stabilizer subgroup of $x$, and suppose that $G_{x} \leq Z$. Suppose $C>0, \epsilon \in(0,1 / 6)$, and if $x \in X$ let us set $m:=\pi(x)$. Then, uniformly in $\mathbf{v}_{1}, \mathbf{v}_{2} \in \mathfrak{g}_{M}\left(m_{x}\right)^{\perp_{h}}$ satisfying $\left\|\mathbf{v}_{j}\right\| \leq C k^{\epsilon}$, belonging to subspace transverse to the orbit through $x$ and with $\theta_{1}-\theta_{2}=0$; we have for $k \rightarrow+\infty$ an asymptotic expansion

$$
\begin{aligned}
\Pi_{k \nu} & \left(x+\frac{1}{\sqrt{k}} v_{1}, x+\frac{1}{\sqrt{k}} v_{2}\right) \\
& \sim \frac{1}{2 \lambda\left(m_{x}\right)}\left(\frac{k \nu}{2 \pi \lambda\left(m_{x}\right)}\right)^{d} \cdot \sum_{g \in G_{x}} f_{1-k}(g) \cdot e^{u_{0} \psi_{2}\left(\mathbf{v}_{1}^{(g)}, \mathbf{v}_{2}\right)} \\
& \cdot\left[1+\sum_{j \geq 1} k^{-j / 2} a_{j}\left(\nu, m, \mathbf{v}_{1}, \mathbf{v}_{2}\right)\right]
\end{aligned}
$$

where $a_{j}(\nu, m, \cdot, \cdot)$ is a polynomial of degree $\leq 3 j$ and parity $(-1)^{j}$.
We can apply Theorems 1.2.3 and 1.2.2 to estimate the dimension of $H(X)_{\nu}$ when $k \rightarrow+\infty$. Let us make this explicit in the case where $\tilde{\mu}$ is generically free.

Corollary 1.2.1. Le us assume that $\mathbf{0} \notin \Phi_{G}(M)$ and that $\tilde{\mu}$ is free on $X$, we have

$$
\lim _{k \rightarrow+\infty}\left[\left(\frac{2 \pi}{k \nu}\right)^{d} \cdot \operatorname{dim}\left(H_{k \nu}(X)\right)\right]=\frac{1}{2} \int_{M} \lambda\left(m_{x}\right)^{-d-1} \mathrm{dV}_{M}(m)
$$

Proof. For $k=1,2, \ldots$ let us define $f_{k} \in \mathcal{C}^{\infty}(M)$ by setting $f_{k}(m):=$ $k^{-d} \Pi_{k \nu}(x, x)$ if $m=\pi(x)$. Given Theorem 1.2.2, $f_{k} \leq C$ for some constant $C>0$ and $f_{k} \rightarrow(\nu / 2 \pi)^{d} \Phi^{-(d+1)} / 2$ for $k \rightarrow+\infty$. By the dominate convergence theorem we can conclude.

### 1.3 The case $G=U(2)$

First, let us set some notations. In this section we will denote by $G$ the group $U(2)$, its Lie algebra $\mathfrak{g}$ is the space of skew-Hermitian matrices. There is a natural $G$-invariant Euclidean scalar product on $\mathfrak{g}$, given by $\left\langle\xi_{1}, \xi_{2}\right\rangle=$ $\operatorname{trace}\left(\xi_{1} \bar{\xi}_{2}^{t}\right)$. By means of $\langle\cdot, \cdot\rangle$, we can equivariantly identify $\mathfrak{g}$ and $\mathfrak{g}^{\vee}$. If $\boldsymbol{\nu} \in \mathbb{Z}^{2}, \nu_{1}>\nu_{2}$, let us set

$$
D_{\nu}:=\left(\begin{array}{cc}
\nu_{1} & 0 \\
0 & \nu_{2}
\end{array}\right) .
$$

Let us introduce the following loci:

1. $\mathcal{O}_{\nu}$ is the (co)adjoint orbit of $i D_{\nu}$;
2. $\mathcal{C}\left(\mathcal{O}_{\nu}\right):=\mathbb{R}_{+} \cdot \mathcal{O}_{\nu}$ is the cone over $\mathcal{O}_{\nu}$;
3. in $M$ and $X$, respectively, we have the inverse images

$$
M_{\mathcal{O}_{\nu}}^{G}:=\Phi_{G}^{-1}\left(\mathcal{C}\left(\mathcal{O}_{\nu}\right)\right), \quad X_{\mathcal{O}_{\nu}}^{G}:=\pi^{-1}\left(M_{\mathcal{O}_{\nu}}^{G}\right)
$$

Finally, let us define the following smooth functions

$$
m \in M_{\mathcal{O}_{\nu}}^{G} \mapsto h_{m} T \in G / T, \quad m \in M_{\mathcal{O}_{\nu}}^{G} \mapsto \lambda_{\nu}(m) \in(0,+\infty)
$$

by the equality

$$
\begin{equation*}
\Phi_{G}(m)=i \lambda_{\boldsymbol{\nu}}(m) h_{m} D_{\boldsymbol{\nu}} h_{m}^{-1} \tag{1.6}
\end{equation*}
$$

Our first result is the following.

Theorem 1.3.1. Assume that $\mathbf{0} \notin \Phi_{G}(M)$, and $\Phi_{G}$ is transverse to $\mathcal{C}\left(\mathcal{O}_{\nu}\right)$. Let us define $G \times G$-invariant subset of $X \times X$

$$
\mathcal{Z}_{\nu}:=\left\{(x, y) \in X_{\mathcal{O}_{\nu}}^{G} \times X_{\mathcal{O}_{\nu}}^{G}: y \in G \cdot x\right\}
$$

Then, uniformly on compact subset of $(X \times X) \backslash \mathcal{Z}_{\boldsymbol{\nu}}$, we have

$$
\Pi_{k \nu}(x, y)=O\left(k^{-\infty}\right)
$$

Corollary 1.3.1. Uniformly on compact subset of $X \backslash X_{\mathcal{O}_{\nu}}^{G}$, we have

$$
\Pi_{k \nu}(x, x)=O\left(k^{-\infty}\right) \quad \text { for } \quad k \rightarrow+\infty
$$

The hypothesis of Theorem 1.3.1 imply that $M_{\mathcal{O}_{\nu}}^{G}$ is a compact and smooth real hypersurface of $M$. Our next aim is to elucidate the geometry of $M_{\mathcal{O}_{\nu}}^{G}$, we need to introduce some further loci related to the action.

Definition 1.3.1. Let

$$
M_{\nu}^{G}:=\Phi_{G}^{-1}\left(i \mathbb{R}_{+} \cdot D_{\nu}\right), \quad X_{\nu}^{G}:=\pi^{-1}\left(M_{\nu}^{G}\right)
$$

Remark 2. Under the assumption of Theorem 1.3.1, $M_{\nu}^{G}$ is a compact submanifold of $M$, of real codimension 3. Clearly, $M_{\mathcal{O}_{\nu}}^{G}=G \cdot M_{\nu}^{G}$, i.e. the G-saturation of $M_{\nu}^{G}$, by the equivariance of $\Phi_{G}$.

There is a natural map diag : $\mathfrak{g}^{\vee} \rightarrow \mathfrak{t}^{\vee}$ given by

$$
\operatorname{diag}: i\left(\begin{array}{ll}
a & z \\
\bar{z} & b
\end{array}\right) \mapsto i\binom{a}{b}
$$

where $a, b \in \mathbb{R}$ and $z \in \mathbb{C}$. The action of $T$ on $M$ induced by restriction of $\mu$ is also Hamiltonian, with moment map given by

$$
\Phi_{T}:=\operatorname{diag} \circ \Phi_{G}: M \rightarrow \mathfrak{t}
$$

Let us introduce the loci

$$
M_{\nu}^{T}:=\Phi_{T}^{-1}\left(\mathbb{R}_{+} \cdot i \boldsymbol{\nu}\right), \quad X_{\nu}^{T}:=\pi^{-1}\left(M_{\nu}^{T}\right)
$$

Let us assume that $\mathbf{0} \notin \Phi_{T}(M)$ and that $\Phi_{T}$ is transverse to $\mathbb{R}_{+} \cdot i \boldsymbol{\nu}$; then $M_{\nu}^{T}$ is a compact smooth real hypersurface of $M$. Since $M_{\nu}^{G} \subset M_{\nu}^{T}$, we have $M_{\mathcal{O}_{\nu}}^{G} \subset G \cdot M_{\nu}^{T}$.

In $\S 4.5 .1$ we shall construct a vector field $\Upsilon=\Upsilon_{\mu, \nu}$ tangent to $M$ along $M_{\mathcal{O}_{\nu}}^{G}$, naturally associated to the action and the weight, which is nowhere vanishing and everywhere normal to $M_{\mathcal{O}_{\nu}}^{G}$. Explicitly, $\Upsilon=\Upsilon_{\mu, \nu}$ along $M_{\mathcal{O}_{\nu}}^{G}$ is

$$
\Upsilon(m):=J_{m}\left(\boldsymbol{\rho}(m)_{M}(m)\right) \quad\left(m \in M_{\mathcal{O}_{\nu}}^{G}\right),
$$

where

$$
\boldsymbol{\rho}(m):=i h_{m} D_{\boldsymbol{\nu}_{\perp}} h_{m}^{-1}
$$

and $\boldsymbol{\nu}_{\perp}:=\left(-\nu_{2}, \nu_{1}\right)^{t}$.
Theorem 1.3.2. Let us assume that:

1. $\Phi_{G}$ and $\Phi_{T}$ are both transverse to $\mathbb{R}_{+} \cdot i D_{\boldsymbol{\nu}}$;
2. $\mathbf{0} \notin \Phi_{T}(M)$ (hence also $\mathbf{0} \notin \Phi_{G}(M)$ );
3. $M_{\nu}^{G} \neq \emptyset$ (equivalently $M_{\mathcal{O}_{\nu}}^{G} \neq \emptyset$ );
4. $\nu_{1}+\nu_{2} \neq 0$.

Then

1. $M_{\mathcal{O}_{\nu}}^{G}$ is a connected and orientable smooth hypersurface in $M$, and separates $M$ in two connected components: the "outside" $A:=M \backslash G \cdot M_{\nu}^{T}$ and the "inside" $B:=G \cdot M_{\nu}^{T} \backslash M_{\mathcal{O}_{\nu}}^{G}$;
2. the normal bundle to $M_{\mathcal{O}_{\nu}}^{G}$ in $M$ is the real line sub-bundle of $\left.T M\right|_{M_{\mathcal{O}_{\nu}}}$ spanned by $\Upsilon$;
3. $\Upsilon$ is "outer" oriented if $\nu_{1}+\nu_{2}>0$ and "inner" oriented if $\nu_{1}+\nu_{2}<0$;
4. $M_{\mathcal{O}_{\nu}}^{G} \cap M_{\nu}^{T}=M_{\nu}^{G}$, and the two hypersurcafes meet tangentially along $M_{\nu}^{G}$.

Let us clarify the meaning of the partition $M=A \dot{\cup} M_{\mathcal{O}_{\nu}}^{G} \dot{\cup} B$. Clearly, $G \cdot M_{\nu}^{T}=\bar{B}, A=\left(G \cdot M_{\nu}^{T}\right)^{c}$. For any $m \in M$, let $\mathcal{O}_{\Phi(m)}:=\Phi_{G}(G \cdot m)$ be the coadjoint orbit of $\Phi_{G}(m)$. By the Schur-Horn Theorem, see [Hor54], $\operatorname{diag}\left(\mathcal{O}_{\Phi(m)}\right)$ is the segment $J_{m}$ joining $i\left(\lambda_{1}, \lambda_{2}\right)^{t}$ with $i\left(\lambda_{2}, \lambda_{1}\right)^{t}$. Then:

1. $m \in A$ if and only if the orthogonal projection of the orbit $\mathcal{O}_{\Phi(m)}$ in $\mathfrak{t}$, $\operatorname{diag}\left(\mathcal{O}_{\Phi(m)}\right)$, is disjoint from $i \mathbb{R}_{+} \cdot \boldsymbol{\nu}$;
2. $m \in M_{\mathcal{O}_{\nu}}^{G}$ if and only if $\operatorname{diag}\left(\mathcal{O}_{\Phi(m)}\right) \cap\left(i \mathbb{R}_{+} \cdot \boldsymbol{\nu}\right)$ is an endpoint of $J_{m}$;
3. $m \in B$ if and only if $\operatorname{diag}\left(\mathcal{O}_{\Phi(m)}\right) \cap\left(i \mathbb{R}_{+} \cdot \boldsymbol{\nu}\right)$ is an interior point of $J_{m}$.

The next step will be to provide some more precise quantitative information on the rate of decay of $\Pi_{k \nu}(\cdot, \cdot)$ on the complement of $\mathcal{Z}_{\nu}$. Namely, we shall show that $\Pi_{k \nu}(x, y)$ is still rapidly decreasing when either $y \rightarrow G \cdot x$ at a sufficiently slow rate, or when at least one of $x$ and $y$ belongs to the "outer" component $A$, and converges to $X_{\mathcal{O}_{\nu}}^{G}$ sufficiently slowly.

Let us consider on $X$ the Riemannian structure which is uniquely determined by the following conditions:

1. (1.2) is an orthogonal direct sum;
2. $\pi: X \rightarrow M$ is a Riemannian submersion;
3. the $S^{1}$-orbits have unit length.

The corresponding density is $\mathrm{dV}_{X}$. Let $\operatorname{dist}_{X}: X \times X \rightarrow[0,+\infty)$ denote the associated distance function.

Theorem 1.3.3. In the situation of Theorem 1.3.1, assume in addition that $G$ acts freely on $X_{\mathcal{O}_{\nu}}^{G}$. For any fixed $C, \epsilon>0$, we have $\Pi_{k \nu}(x, y)=O\left(k^{-\infty}\right)$ uniformly for

$$
\max \left\{\operatorname{dist}_{X}(x, G \cdot y), \operatorname{dist}_{X}\left(x, G \cdot X_{\nu}^{T}\right)\right\} \geq C k^{\epsilon-1 / 2}
$$

Remark 3. Theorem 1.3.3 is built on the general theory of Guillemin and Sternberg; we have established that the $U(2)$-equivariant asymtptotics of the Szegö kernels are rapidly decreasing for points not lying in $X_{\mathcal{O}_{\nu}}^{G}$. However, the precise structure of this rapid decrease along the boundary is not clear at the moment, since the usual scaling techniques based on [SZ02] do not seem to cope with "inner" directions.

In Theorem 1.3.4 below we shall consider the diagonal asymptotics behaviour of $\Pi_{k \nu}$ along $X_{\mathcal{O}_{\nu}}^{G}$ assuming the action is free on it. Before giving the statement, some further notation is needed.

Definition 1.3.2. Given $\boldsymbol{\nu} \in \mathbb{Z}^{2}$, we shall denote by $\boldsymbol{\nu}_{M} \in \mathfrak{X}(M)$ and $\boldsymbol{\nu}_{X} \in \mathfrak{X}(X)$ the vector fields induced by $i D_{\boldsymbol{\nu}}$; similarly, for every $g T \in$ $G / T, \operatorname{Ad}_{g}(\boldsymbol{\nu})_{M}$ and $\operatorname{Ad}_{g}(\boldsymbol{\nu})_{X}$ will be the vector field induced by $\operatorname{Ad}_{g}\left(i D_{\boldsymbol{\nu}}\right)$ respectively on $M$ and $X$.

Definition 1.3.3. Let $\|\cdot\|_{m}: T_{m} M \rightarrow \mathbb{R}$ and $\|\cdot\|_{x}: T_{x} X \rightarrow \mathbb{R}$ be the norm functions. If $\boldsymbol{\nu}_{1}, \boldsymbol{\nu}_{2} \in \mathbb{Z}^{2}, \nu_{1}>\nu_{2}$, let us set $\boldsymbol{\nu}_{\perp}:=\left(-\nu_{2}, \nu_{1}\right)$. Let us define the smooth function $\mathcal{D}_{\nu}: M_{\mathcal{O}_{\nu}}^{G} \rightarrow(0,+\infty)$ by posing

$$
\mathcal{D}_{\nu}(m):=\frac{\|\boldsymbol{\nu}\|}{\left\|\operatorname{Ad}_{h_{m}}\left(\boldsymbol{\nu}_{\perp}\right)_{M}(m)\right\|_{m}}
$$

For every $m \in M_{\mathcal{O}_{\nu}}^{G}$ and $g T \in G / T$ we have $\left\langle\Phi_{G}(m), \operatorname{Ad}_{g}\left(i D_{\nu_{\perp}}\right)\right\rangle=0$. Let us set $\pi(x)=m_{x}$. Hence, by (1.4), the function $\mathcal{D}_{\nu}$ is well-defined since

$$
\left\|\operatorname{Ad}_{h_{m}}\left(\boldsymbol{\nu}_{\perp}\right)_{M}\left(m_{x}\right)\right\|_{m_{x}}=\left\|\operatorname{Ad}_{h_{m_{x}}}\left(\boldsymbol{\nu}_{\perp}\right)_{X}(x)\right\|_{x}>0 .
$$

Let us record one more piece of notation. If $V_{3}$ is the area of the unit sphere $S^{3} \subset \mathbb{R}^{4}$, let us set

$$
D_{G / T}:=2 \pi V_{3}^{-1}
$$

Theorem 1.3.4. Under the same hypothesis as in Theorem 1.3.2, let us assume in addition that $G$ atcs freely on $X_{\mathcal{O}_{\nu}}^{G}$. Then uniformly in $x \in X_{\mathcal{O}_{\nu}}^{G}$ we have for $k \rightarrow+\infty$ an asymptotics expansion of the form

$$
\begin{gathered}
\Pi_{k \nu}(x, x) \sim \frac{D_{G / T}}{\sqrt{2}} \frac{1}{\left\|\Phi_{G}\left(m_{x}\right)\right\|^{d+1 / 2}\left(\frac{k\|\boldsymbol{\nu}\|}{\pi}\right)^{d-1 / 2} \cdot \mathcal{D}_{\boldsymbol{\nu}}\left(m_{x}\right)} \\
{\left[1+\sum_{j \geq 1} k^{-j / 4} a_{j}\left(\boldsymbol{\nu}, m_{x}\right)\right] .}
\end{gathered}
$$

Let us also discuss how the expansion generalizes to the case where the action on $X_{\mathcal{O}_{\nu}}^{G}$ is only locally free. Let $x \in X_{\mathcal{O}_{\nu}}^{G}$, then the stabilizer $G_{x}$ is a discrete, hence finite, subgroup of $G$. By the equivariance of $\Phi_{G}$, we have that $g$ commutes with $\Phi_{G}\left(m_{x}\right)$ for each $g \in G_{m_{x}}$. Thus, let $h_{m_{x}}$ as in (1.6), all the matrices that commute with $\Phi_{G}\left(m_{x}\right)$ have the form

$$
h_{m_{x}}\left(\begin{array}{cc}
e^{i \alpha} & 0 \\
0 & e^{i \beta}
\end{array}\right) h_{m_{x}}^{-1},
$$

for some $e^{i \alpha}$, $e^{i \beta} \in S^{1}$. Since $G_{x} \subseteq G_{m_{x}}$, there exists a well-defined finite subgroup $R_{x}$ of $T$ such that $G_{x}=h_{m} R_{x} h_{m}^{-1}$.

Let $Z$ the subgroup of scalar matrices, $Z \cong S^{1}$, and denote with $Z_{x}:=$ $Z \cap G_{x}$. Since $Z_{x}$ is finite subgroup of $S^{1}$, it is a cyclic group, i.e. there exists $b_{x} \in \mathbb{N}$ such that

$$
\beta_{x}=e^{\frac{2 \pi}{b_{x}} i} I_{2}
$$

is the generator of $Z_{x},\left|Z_{x}\right|=b_{x}$. Let us define the function

$$
\Gamma_{k}(x)=\sum_{g \in Z_{x}} g^{k}=\sum_{j=0}^{b_{x}-1}\left(\beta_{x}^{k}\right)^{j} .
$$

Furthermore $G_{x} \backslash Z_{x}$ has even cardinality $2 a_{x}$, and perhaps after renumbering its elements can be arranged in pairs $\left(g_{j}, g_{j}^{-1}\right), j=1, \ldots, a_{x}$. Thus for each $j=1,2, \ldots, a_{x}$ the corresponding $g_{j} \in G_{x} \backslash Z_{x}$ has the form

$$
g_{j}=h_{m}\left(\begin{array}{cc}
e^{i \alpha_{j}} & 0 \\
0 & e^{i \beta_{j}}
\end{array}\right) h_{m}^{-1},
$$

where $\alpha_{j} \neq \beta_{j}$. In $\S 4.8$ we shall define a complex-valued smooth function $\Theta_{j}(x, \boldsymbol{\nu})$, depending on $\alpha_{j}, \beta_{j}, \boldsymbol{\nu}$ and $x$; for ease of exposition we postpone its definition and the proof of the following Theorem 1.3.5 in section §4.8.

Theorem 1.3.5. Assume that the same hypothesis as in Theorem 1.3.2 hold. Then uniformly in $x \in X_{\mathcal{O}_{\nu}}^{G}$ we have for $k \rightarrow+\infty$ that the leading term of the asymptotics expansion of $\Pi_{k \nu}(x, x)$ is

$$
\begin{aligned}
& \frac{D_{G / T}}{\left\|\Phi_{G}\left(m_{x}\right)\right\|^{d+1 / 2}} \cdot\left(\frac{k\|\boldsymbol{\nu}\|}{\pi}\right)^{d-1 / 2}\left\{\frac{1}{\sqrt{2}} \mathcal{D}_{\boldsymbol{\nu}}\left(m_{x}\right) \Gamma_{1-\left(\nu_{1}+\nu_{2}\right) k}(x)\right. \\
& \left.\quad+2 \sqrt{2\left\|\Phi_{G}\left(m_{x}\right)\right\|\|\boldsymbol{\nu}\|}\left(\nu_{1}-\nu_{2}\right) \cdot \sum_{j=0}^{a_{x}} \Re\left[\frac{\left(e^{i \alpha_{j}}-e^{i \beta_{j}}\right) e^{-i k\left(\nu_{1} \alpha_{j}+\nu_{2} \beta_{j}\right)}}{\sqrt{\Theta_{j}(x, \boldsymbol{\nu})}}\right]\right\} .
\end{aligned}
$$

From now on, to simplify our exposition, we shall make the stronger assumption that $\tilde{\mu}$ is actually free along $X_{\mathcal{O}_{\nu}}^{G}$.

We can refine the previous asymptotic expansion at a fixed diagonal point $(x, x) \in X_{\mathcal{O}_{\nu}}^{G} \times X_{\mathcal{O}_{\nu}}^{G}$ to an asymptotic expansion for near-diagonal rescaled displacement; however, for the sake of simplicity we shall restrict the directions of the displacements.

Definition 1.3.4. If $m \in M$, let $\mathfrak{g}_{M}(m) \subset T_{m} M$ be the image of the linear evaluation map $\operatorname{val}_{m}: \mathfrak{g} \rightarrow T_{m} M, \xi \mapsto \xi_{M}(m)$; also, let $\mathfrak{g}_{M}(m)^{\perp_{\omega}} \subseteq T_{m} M$ be its symplectic orthocomplement with respect to $\omega_{m}$, and let $\mathfrak{g}_{M}(m)^{\perp_{g}} \subseteq T_{m} M$ be its Riemannian orthocomplement with respect to $g_{m}$. Hence,

$$
\mathfrak{g}_{M}(m)^{\perp_{h}}:=\mathfrak{g}_{M}(m)^{\perp_{\omega}} \cap \mathfrak{g}_{M}(m)^{\perp_{g}} \subseteq T_{m} M
$$

is the Hermitian orthocomplement of the complex subspace generated by $\mathfrak{g}_{M}(m)$ with respect to $h_{m}:=g_{m}-i \omega_{m}$.

Definition 1.3.5. If $\mathbf{v}_{1}, \mathbf{v}_{2} \in T_{m} M$, following [SZ02] let us set

$$
\begin{equation*}
\psi_{2}\left(\mathbf{v}_{1}, \mathbf{v}_{2}\right):=-i \omega_{m}\left(\mathbf{v}_{1}, \mathbf{v}_{2}\right)-\frac{1}{2}\left\|\mathbf{v}_{1}-\mathbf{v}_{2}\right\|_{m}^{2} \tag{1.7}
\end{equation*}
$$

Here $\|\mathbf{v}\|_{m}:=g_{m}(\mathbf{v}, \mathbf{v})^{1 / 2}$. The same invariant can be introduced in any Hermitian vector space. Given the choice of a system of Heisenberg local coordinates centred at $x \in X$ ([SZ02]), there is built-in unitary isomorphism $T_{m} M \cong \mathbb{C}^{d} ;$ with this implicit, (1.7) will be used with $\mathbf{v}_{j} \in \mathbb{C}^{d}$.

The choice of Heisenberg local coordinates centred at $x \in X$ gives a meaning to the expression $x+(\theta, \mathbf{v})$ for $(\theta, \mathbf{v}) \in(\pi, \pi) \times \mathbb{R}^{2 d}$ with $\|\mathbf{v}\|$ of sufficiently small norm. When $\theta=0$, we shall write $x+\mathbf{v}$.

Theorem 1.3.6. Let us assume the same hypothesis as in Theorem 1.3.4. Suppose $C>0, \epsilon \in(0,1 / 6)$, and if $x \in X$ let us set $m_{x}:=\pi(x)$. Then, uniformly in $x \in X_{\mathcal{O}_{\nu}}^{G}$ and $\mathbf{v}_{1}, \mathbf{v}_{2} \in \mathfrak{g}_{M}\left(m_{x}\right)^{\perp_{h}}$ satisfying $\mathbf{v}_{j} \leq C k^{\epsilon}$, we have for $k \rightarrow+\infty$ an asymptotics expansion

$$
\begin{aligned}
\Pi_{k \nu}\left(x+\frac{1}{\sqrt{k}} \mathbf{v}_{1}, x+\frac{1}{\sqrt{k}} \mathbf{v}_{2}\right) \sim & \frac{D_{G / T}}{\sqrt{2}} \frac{e^{\psi_{2}\left(\mathbf{v}_{1}, \mathbf{v}_{2}\right) / \lambda_{\nu}\left(m_{x}\right)}}{\left\|\Phi_{G}\left(m_{x}\right)\right\|^{d+1 / 2}}\left(\frac{k\|\boldsymbol{\nu}\|}{\pi}\right)^{d-1 / 2} \\
& \cdot \mathcal{D}_{\nu}\left(m_{x}\right)\left[1+\sum_{j \geq 1} k^{-j / 4} a_{j}\left(\nu, m_{x}, \mathbf{v}_{1}, \mathbf{v}_{2}\right)\right]
\end{aligned}
$$

where $a_{j}(\nu, m, \cdot, \cdot)$ is a polynomial of degree $\leq\lceil 3 j / 2\rceil$.
Furthermore, we shall provide an integral formula of independent interest for the asymptotics of $\Pi_{k \nu}\left(x^{\prime}, x^{\prime}\right)$ when $x^{\prime} \rightarrow X_{\mathcal{O}_{\nu}}^{G}$ at a "fast" pace from the "outside" (that is, $x^{\prime} \in A$ in the notation of Theorem 1.3.2). While the latter formula is a bit too technical to be described in this introduction, by global integration it leads to a lower bound on $\operatorname{dim} H(X)_{\nu}$ which can be stated in a compact form. With the notation of Theorem 1.3.2, we have

$$
\operatorname{dim} H(X)_{\boldsymbol{\nu}}=\operatorname{dim}_{\text {in }} H(X)_{\boldsymbol{\nu}}+\operatorname{dim}_{\text {out }} H(X)_{\boldsymbol{\nu}}
$$

where

$$
\operatorname{dim}_{\text {out }} H(X)_{\nu}:=\int_{A} \Pi_{\nu}(x, x) \mathrm{dV}_{X}(x)
$$

and similarly for $\operatorname{dim}_{i n} H(X)_{\nu}$, with $A$ replaced by $B$. Hence an asymptotic estimate for $\operatorname{dim}_{\text {out }} H(X)_{\nu}$ when $k \rightarrow+\infty$ implies an asymptotic lower bound
for $\operatorname{dim} H(X)_{k \nu}$. In Theorem 1.3.7 below, we shall show that $\operatorname{dim}_{o u t} H(X)_{k \nu}$ is given by an asymptotic expansion of descending fractional powers of $k$, the leading power being $k^{d-1}$.

Theorem 1.3.7. Under the assumption of Theorem 1.3.4, $\operatorname{dim}_{\text {out }} H(X)_{\nu}$ is given by an asymptotic expansion in descending powers of $k^{1 / 4}$ as $k \rightarrow+\infty$, with leading order term

$$
\frac{1}{4} D_{G / T}\left(\frac{k\|\boldsymbol{\nu}\|}{\pi}\right)^{d-1} \int_{M_{O_{\boldsymbol{\nu}}}} \frac{1}{\left\|\Phi_{G}(m)\right\|^{d}} \cdot \mathcal{D}_{\boldsymbol{\nu}}(m) \mathrm{dV}_{M_{O_{\boldsymbol{\nu}}}^{G}}(m)
$$

### 1.4 Recalls on prior literature

The asymptotic expansions of the Szegö kernel and its variants were studied in a lot of papers and they are grounded on a well-known foundational result due to Boutet de Monvel and Sjöstrand [BdMS76], who proved that $\Pi$ is a Fourier integral operator with complex phase. Beginning with the papers of Zelditch on a theorem of Tian, [Zel98], the FIO construction for these kernel functions has proved extremely useful in the study of the behavior, as the power tends to infinity, of the space of sections of powers of a positive line bundle over an algebraic variety.

More explicitly, let $(A, h)$ be an Hermitian line bundle over a compact $d$-dimensional Kähler manifold $(M, \omega)$. Let $\left(s_{0}^{k}, \ldots, s_{d_{k}}^{k}\right)$ be any orthonormal basis of $H^{0}\left(M, A^{\otimes k}\right)$ (where $d_{k}+1$ is the dimension of $H^{0}\left(M, A^{\otimes k}\right)$ ), with respect to the inner product

$$
\left\langle s_{1}, s_{2}\right\rangle_{h_{k}}=\int_{M} h_{k}\left(s_{1}(m), s_{2}(m)\right) \mathrm{dV}_{M}(m)
$$

Catlin and Zelditch, indipendently, sharpened a theorem conjectured by Yau, solved by Tian, proving that there exists a complete asymptotic expansion:

$$
\begin{equation*}
\sum_{j=0}^{d_{k}}\left\|s_{j}^{k}(m)\right\|_{h_{k}}^{2}=a_{0} k^{d}+a_{1}(m) k^{d-1}+a_{2}(m) k^{d-2}+\ldots \tag{1.8}
\end{equation*}
$$

for certain smooth coefficients $a_{j}(m)$ with $a_{0}=\pi^{-d}$. More recently, Z. Lu [Lu00] proved that the lower terms of the asymptotic expansion (1.8) is a polynomial of the curvature and its covariant derivatives at $m$ of the metric
$g$, associated to the form $\omega$; and he computed the first three coefficients of this expansion (see also [Loi04] and [Loi05] for the computations of the coefficients $a_{j}$ 's).

The expansion (1.8) is sometimes called the TYZ (Tian-Yau-Zelditch) expansion (also Catlin in [Cat99] published the same result as Zelditch in the same year). Since the CR structure of the circle bundle $X$ is $S^{1}$-invariant (see the next chapter for definitions and details), there is a naturally induced unitary representation of $S^{1}$ on $H(X)$; therefore, $H(X)$ splits unitarily and equivariantly as a direct sum of isotypical components $H_{k}(X), k \in \mathbb{Z}$. We have $H_{k}(X)=0$ if $k<0$, and for $k \geq 0$ there is a natural unitary isomorphism between $H_{k}(X)$ and $H^{0}\left(M, A^{\otimes \bar{k}}\right)$. Thus the expansion (1.8) coincides with the asymptotics of the kernel of the projector $\Pi_{k}: L^{2}(X) \rightarrow H_{k}(X)$ along the diagonal.

In the situation described above the standard $S^{1}$ action on $X$ is trivial on $M$ and thus it is Hamiltonian with respect to $2 \omega$ with constant moment map. Thus we can generalize the situation described in the previous paragraph as follow. Suppose that we have another Hamiltonian holomorphic action $\mu$ of $S^{1}$ on $M$ with moment map $\Phi: M \rightarrow \mathbb{R}$; suppose furthermore that it can be linearized to a holomorphic action $\tilde{\mu}$ on $A$. Then $S^{1}$ acts on $X$ as a group of contactomorphisms under the naturally induced action. Since the lifted action $\tilde{\mu}$ preserves the contact form $\alpha$ and it is a lifting of the holomorphic action $\mu$, it leaves $H(X)$ invariant; therefore it determines a unitary action of $S^{1}$ on $H(X)$. Thus $H(X)$ equivariantly and unitarily decomposes into the Hilbert direct sum of its isotypes $H_{k}^{\mu}(X)$. In general, however, the isotypical components in point don't correspond to subspaces of holomorphic sections of some higher tensor power of the polarizing line bundle.

Nonetheless, if $\Phi>0$ then $H_{k}^{\mu}(X)$ is finite-dimensional for any $k \in \mathbb{Z}$, and is the null space if $k<0$; in particular, the orthogonal projector $\Pi_{k}^{\mu}(X)$ : $L^{2}(X) \rightarrow H_{k}^{\mu}(X)$ is a smoothing operator. Uniformly on $x \in X$, there is an asymptotic expansion of the form

$$
\Pi_{k}^{\mu}(x, x) \sim\left(\frac{k}{\pi}\right)^{d} \cdot\left[\Phi(m)^{-(d+1)}+\sum_{j \geq 1} k^{-j} a_{j}^{\mu}\left(m_{x}\right)\right]
$$

and the first term $a_{1}^{\mu}$ is computed in [Pao15]; in particular if $\Phi=1$, one recovers Lu's subprincipal term [Lu00].

In the $n$-dimensional toric case, this theme has been studied in [Pao12], [Pao15] and [Cam16]. One has a decomposition of the Hardy space into
isotypes $H_{k \nu}(X), \boldsymbol{\nu} \in \mathbb{Z}^{n}$, and one is led to investigate the behavior of the kernel of the projector $\Pi_{k \nu}: L^{2}(X) \rightarrow H_{k \nu}(X)$ as $k$ goes to infinity. Suppose that $\mathbf{0} \notin \Phi(M)$; then

1. If $\Phi(m) \notin \mathbb{R}_{+} \cdot \boldsymbol{\nu}$, then $\Pi_{k \nu}=O\left(k^{-\infty}\right)$.
2. Assume that $\Phi$ is transversal to $\mathbb{R}_{+} \cdot \boldsymbol{\nu}$. Then for every $m=\pi(x) \in$ $\Phi^{-1}\left(\mathbb{R}_{+} \cdot \boldsymbol{\nu}\right)$ (where $\pi$ is the projection $\left.\pi: A \rightarrow M\right)$ as $k \rightarrow \infty$ we have

$$
\begin{aligned}
\Pi_{k \nu}(x, x) \sim & \frac{1}{(\sqrt{2} \pi)^{n-1}}\left(\|\boldsymbol{\nu}\| \cdot \frac{k}{\pi}\right)^{d+(1-n) / 2} \cdot \sum_{g \in T_{m}} \chi_{\boldsymbol{\nu}}(g)^{k} \\
& \cdot \frac{1}{\mathcal{D}(m)}\left(\frac{1}{\|\Phi(m)\|}\right)^{d+1+(1-n) / 2} \cdot\left(1+\sum_{l \geq 1} B_{l}\left(m_{x}\right) k^{-l}\right)
\end{aligned}
$$

where $T_{m}$ is the stabilizer of $m, \chi_{\nu}$ is the character pertaining to $\boldsymbol{\nu}$, $\mathcal{D}: M \rightarrow \mathbb{R}$ is a distortion function and $B_{l}$ 's are smooth functions on $\Phi^{-1}\left(\mathbb{R}_{+} \cdot \boldsymbol{\nu}\right)$ (we refer to [Pao12] for definitions).

In general, if the standard circle action is replaced by a geometrically induced representation of a generic compact Lie group $G$, then, under appropriate hypothesis, the isotypes are finite dimensional and one can ask whether there are some extension of the above asymptotic expansions. In this thesis, we turn to non-Abelian actions, and consider specifically the cases of $G=S U(2)$ and $G=U(2)$.

The analysis of Szegö kernels have other applications in complex and symplectic geometry. By studying the off-diagonal of $\Pi_{k}$, Shiffman and Zelditch obtain an analytic proof of the Kodaira embedding theorem (see [SZ02]). Explicitly, when $x \in X$ tends to the orbit through $y$ at a suitable fast pace, there is an asymptotic expansion which captures an exponential decrease of $\Pi_{k}(x, y)$ away from $(\pi \times \pi)^{-1}\left(\Delta_{M}\right)$, in a family of shrinking neighborhoods of the latter ( $\Delta_{M}$ is the diagonal in $M \times M$ ). To express this, following [BSZ00] and [SZ02] let us define $\psi_{2}: \mathbb{C}^{d} \times \mathbb{C}^{d} \rightarrow \mathbb{C}$ by

$$
\psi_{2}\left(\mathbf{v}_{1}, \mathbf{v}_{2}\right):=-\imath \omega\left(\mathbf{v}_{1}, \mathbf{v}_{2}\right)-\frac{1}{2}\left\|\mathbf{v}_{1}-\mathbf{v}_{2}\right\|^{2}
$$

in the so-called Heisenberg local coordinates (henceforth, HLC's) centered at
$x \in X$, where $v_{l}=\left(\theta, \mathbf{v}_{l}\right) \in \mathbb{R} \otimes \mathbb{C}^{d} \cong T_{x} X$. We have

$$
\begin{align*}
\Pi_{k}\left(x+\left(\theta_{1}, \frac{\mathbf{v}_{1}}{\sqrt{k}}\right), x+\left(\theta_{2}, \frac{\mathbf{v}_{2}}{\sqrt{k}}\right)\right) \sim & \left(\frac{k}{\pi}\right)^{d} \cdot e^{\imath k\left(\theta_{1}-\theta_{2}\right)+\psi_{2}\left(\mathbf{v}_{1}, \mathbf{v}_{2}\right)}  \tag{1.9}\\
& \cdot\left[1+\sum_{j \geq 1} k^{-j / 2} \cdot R_{j}\left(m_{x} ; \mathbf{v}_{1}, \mathbf{v}_{2}\right)\right]
\end{align*}
$$

Here $R_{j}\left(m_{x} ; \cdot, \cdot\right)$ is a polynomial of degree $\leq 3 j$. For fixed $C>0$ and $\epsilon \in(0,1 / 6)$, the asymptotic expansion (1.9) holds uniformly for $\left\|\mathbf{v}_{j}\right\| \leq C k^{\epsilon}$.

For a general torus action we have, under some suitable hypothesis and uniformly in $m_{x} \in \Phi^{-1}\left(\mathbb{R}_{+} \cdot \boldsymbol{\nu}\right)$, the following asymptotic expansion as $k \rightarrow$ $+\infty$

$$
\begin{align*}
\Pi_{k \nu} & \left(x+\frac{\mathbf{v}_{1}}{\sqrt{k}}, x+\frac{\mathbf{v}_{2}}{\sqrt{k}}\right) \sim  \tag{1.10}\\
& \frac{1}{(\sqrt{2} \pi)^{n-1}}\left(\|\boldsymbol{\nu}\| \cdot \frac{k}{\pi}\right)^{d+(1-n) / 2} \frac{1}{\mathcal{D}(m)}\left(\frac{1}{\|\Phi(m)\|}\right)^{(d+1)+(1-n) / 2} \\
& \cdot\left(\sum_{t \in T_{m}} \chi_{\boldsymbol{\nu}}(t)^{k} \cdot e^{H_{m}\left(\mathrm{~d}_{m} \tilde{\mu}_{t}-1\left(\mathbf{v}_{1}\right), \mathbf{v}_{2}\right)}\right) \\
& \cdot\left[1+\sum_{j \geq 1} k^{-j / 2} \cdot R_{j}\left(m_{x} ; \mathbf{v}_{1}, \mathbf{v}_{2}\right)\right]
\end{align*}
$$

where $\mathbf{v}_{l} \in N_{m} \subseteq T_{m} M$ (here $N$ is the normal bundle of $\Phi^{-1}\left(\mathbb{R}_{+} \cdot \boldsymbol{\nu}\right)$ in $M$ ) and $\left\|\mathbf{v}_{l}\right\| \leq C k^{1 / 9}$. The smooth functions $R_{j}$ 's are polynomial in the $\mathbf{v}_{l}$ 's and $H: T M \oplus T M \rightarrow \mathbb{C}$ is a smooth function (see [Pao12] for the precise definition). We provided similar results for the case of $G=S U(2)$ and $G=U(2)$, see Theorem 1.2.3 and Theorem 1.3.6. While we have restricted the exposition to the complex projective setting, the results of these thesis admit natural generalizations to the almost Kähler context, following the theory of generalized Szegö kernels in [SZ02].

In the paper [SZ99] Shiffman and Zelditch also use the Szegö kernel to show that the zeros of a 'random section' of $H^{0}\left(M, A^{\otimes k}\right)$ become uniformly distributed as $k \rightarrow+\infty$; it is also a key ingredient in the investigations of balanced metrics in Donaldsons terminology [Don01] (see also [AL04]). For variants of and alternative approaches to the general theme of asymptotic expansions see for instance [BU00], [Cha16], [MM07] and [MZ05].

In closing, let us emphasize that while the present analysis and results belong to the general framework of geometric quantization, since the unitary representation on $H(X)$ is a 'quantization' of the Hamiltonin action on $M$, they do not fit into the traditional framework of Berezin-Toeplitz quantization, since we are not working here within a fixed isotype for the structure $S^{1}$-action. Rather, a more appropriate heuristic framework is the one discussed in [GS82c] (and of course [BdMG81]). In this respect, let us note that, at least heuristically, the isotypical components in point should relate to the Riemann Roch numbers of certain symplectic reductions, of which the local asymtptotic expansions above can be seen as an estimate and geometric reinterpretation.

## Chapter 2

## Examples

### 2.1 Examples about $G=S U(2)$

Let us describe some examples where the geometric hypothesis of the previous discussion are verified.

## Example 1

Let $A$ be the hyperplane line bundle on $M=\mathbb{P}^{1}$, then the circle bundle $X \subset A^{\vee}$ may be identified with $S^{3} \subset \mathbb{C}^{2}$, and the projection $\pi: X \rightarrow M$ with the Hopf map. There is a natural action of $G$ on $X$ given by left translation on itself. This action is free and descend to an Hamiltonian action on $M$ given by

$$
\mu_{A}([Z]):=[A Z], \quad\left(Z=\left(z_{0}, z_{1}\right) \in \mathbb{C}^{2} \backslash\{0\}\right)
$$

The moment map is

$$
\Psi([Z]):=\frac{1}{\|Z\|^{2}}\left(\begin{array}{cc}
\frac{1}{2}\left(\left|z_{0}\right|^{2}-|z|_{1}^{2}\right) & z_{0} \bar{z}_{1} \\
\bar{z}_{0} z_{1} & \frac{1}{2}\left(\left|z_{1}\right|^{2}-|z|_{0}^{2}\right)
\end{array}\right)
$$

which is everywhere non-vanishing. Then $\lambda\left(\left[z_{0}: z_{1}\right]\right)=1 / 2$ for any $\left[z_{0}\right.$ : $\left.z_{1}\right] \in \mathbb{P}^{1}$, and the contact action $\tilde{\mu}$ on $S^{3}$ is free, since it may be identified with action of $S U(2)$ on itself by left translations. Furthermore, $H_{k \nu}(X)=$ $H_{k \nu-1}(X)$, where the right hand side is the $(k \cdot \nu-1)$-th isotype for the $S^{1}$-action. With $\nu=1$ the leading order term of the expansion of Theorem 1.2.2 is $(k / \pi)^{d} \cdot e^{\psi_{2}\left(\mathbf{v}_{1}, \mathbf{v}_{2}\right)}$, in agreement with the standard off-diagonal scaling asymptotics for Szegö kernels on $\mathbb{P}^{1}$, [SZ02].

## Example 2

Let us consider the diagonal action of $G$ on $\mathbb{P}^{1} \times \mathbb{P}^{1}$,

$$
\mu_{A}([Z],[W])=([A Z],[A W])
$$

For $r=1,2, \ldots$, consider the symplectic structure $\Omega_{r}:=\omega_{F S} \boxtimes\left(r \omega_{F S}\right)$ on $\mathbb{P}^{1} \times \mathbb{P}^{1}$. Then $\mu$ is Hamiltonian with respect to $2 \Omega_{r}$, with the moment map

$$
\Phi_{r}:([Z],[W]) \mapsto \Psi([Z])+r \Psi([W]) .
$$

If $r \geq 2$, then $\Phi_{r}$ is nowhere vanishing.
On the other hand, $\Omega_{r}$ is the normalized curvature of the positive line bundle $A_{r}:=\mathcal{O}_{\mathbb{P}^{1}}(1) \boxtimes \mathcal{O}_{\mathbb{P}^{1}}(r)$. The unit circle bundle $X_{r}$ associated to $A_{r}$ is the image of $S^{3} \times S^{3}$ under the map

$$
(Z, W) \in S^{3} \times S^{3} \subset \mathbb{C}^{2} \times \mathbb{C}^{2} \mapsto Z \otimes W^{\otimes r} \in \mathbb{C}^{2(r+1)},
$$

and the contact lift of $\mu$ is given by

$$
\tilde{\mu}_{A}\left(Z \otimes W^{\otimes r}\right)=(A Z) \otimes(A W)^{\otimes r}
$$

Let us consider the stabilizer subgroup of $Z \otimes W^{\otimes r}$. We have

$$
\tilde{\mu}_{A}\left(Z \otimes W^{\otimes r}\right)=Z \otimes W^{\otimes r} \Leftrightarrow A Z=\lambda_{1} Z, \quad A W=\lambda_{2} W
$$

for certain $\lambda_{1}, \lambda_{2} \in S^{1}$ with $\lambda_{1} \cdot \lambda_{2}^{r}=1$.
If $Z$ and $W$ are linearly dependent, then $\lambda_{1}=\lambda_{2}$ and $\lambda_{1}^{r+1}=1$. The stabilizer subgroup of $Z \otimes W^{\otimes r}$ is therefore cyclic of order $r+1$. Otherwise, $(Z, W)$ is an eigenbasis of $A$ and $\lambda_{2}=\lambda_{1}^{-1}, \lambda_{1}^{r-1}=1$. Hence, assuming that $Z \wedge W \neq 0$, the stabilizer subgroup of $Z \otimes W^{\otimes r}$ is cyclic of order $r-1$ when $(Z, W)$ is an orthonormal basis of $\mathbb{C}^{2}$, and otherwise it is trivial when $r$ is even and $\left\{ \pm I_{2}\right\}$ when $r$ is odd. Thus $\tilde{\mu}$ is locally free for $r \geq 2$. Furthermore, the action is generically free when $r$ is even, and the stabilizer is generically of order two when $r$ is odd.

Let us now consider how $V_{k \nu}$ appears in

$$
H\left(X_{r}\right)=\bigoplus_{l=0}^{+\infty} H_{l}\left(X_{r}\right), \quad H_{l}\left(X_{r}\right) \cong H^{0}\left(\mathbb{P}^{1} \times \mathbb{P}^{1}, A_{r}^{\otimes l}\right)
$$

Since $A_{r}^{\otimes l}=\mathcal{O}_{\mathbb{P}^{1}}(l) \boxtimes \mathcal{O}_{\mathbb{P}^{1}}(l r)$, by the Künneth formula we have

$$
\begin{aligned}
H^{0}\left(\mathbb{P}^{1} \times \mathbb{P}^{1}, A_{r}^{\otimes l}\right) & \cong H^{0}\left(\mathbb{P}^{1}, \mathcal{O}_{\mathbb{P}^{1}}(l)\right) \otimes H^{0}\left(\mathbb{P}^{1}, \mathcal{O}_{\mathbb{P}^{1}}(l r)\right) \\
& \cong V_{(l+1,0)} \otimes V_{(l r+1,0)}
\end{aligned}
$$

Thus the character of $H^{0}\left(\mathbb{P}^{1} \times \mathbb{P}^{1}, A_{r}^{\otimes l}\right)$ as a $G$-representation is $\chi_{l+1} \cdot \chi_{l r+1}$. We see by a few computations that

$$
\begin{aligned}
\left(\chi_{l+1}\right. & \left.\cdot \chi_{l r+1}\right)\left(\left(\begin{array}{cc}
e^{i \vartheta} & 0 \\
0 & e^{-i \vartheta}
\end{array}\right)\right) \\
& =\left(e^{i l \theta}+e^{i(l-2) \theta}+\cdots+e^{-i l \theta}\right) \cdot \frac{e^{i(l r+1) \theta}-e^{-i(l r+1) \theta}}{e^{i \theta}-e^{-i \theta}} \\
& =\sum_{j=0}^{l} \frac{e^{i(l+l r+1-2 j) \theta}-e^{-i(l+l r+1-2 j) \theta}}{e^{i \theta}-e^{-i \theta}} .
\end{aligned}
$$

Therefore,

$$
H^{0}\left(\mathbb{P}^{1} \times \mathbb{P}^{1}, A_{r}^{\otimes l}\right) \cong \bigoplus_{j=0}^{l} V_{(l+l r+1-2 j, 0)}
$$

We conclude that $V_{k \nu}$ appears at most once in each $H_{l}\left(X_{r}\right)$; it does appear once, in fact, if and only if $k \nu$ and $l(r+1)+1$ have the same parity, and

$$
\begin{equation*}
\frac{k \nu-1}{r-1} \geq l \geq \frac{k \nu-1}{r+1} \tag{2.1}
\end{equation*}
$$

Suppose, for example, that $k \nu$ and $r+1$ are both even. Then $l(r+1)+1$ is odd for any choice of $l$ and we conclude that $H_{k \nu}(X)$ vanishes. Notice that at the general $x \in X_{r}$ we have $G_{x}=\left\{ \pm I_{2}\right\}$, and $\sum_{g \in G_{x}} f_{1-k \nu}(g)=0$. If, on the other hand, $r+1$ is even and $k \nu$ is odd, then there is a copy of $V_{k \nu}$ in $H_{l}\left(X_{r}\right)$ for every integer $l$ satisfying (2.1). Hence the number of copies of $V_{k \nu}$ in $H\left(X_{r}\right)$ is $\sim 2 k \nu /\left(r^{2}-1\right)$, so that the dimension of $H_{k \nu}(X)$ is $\sim 2(k \nu)^{2} /\left(r^{2}-\right.$ 1). For the general $x \in X_{r}$, we have in this case $\sum_{g \in G_{x}} f_{1-k \nu}(g)=2$.

When $r+1$ is odd, on the other hand, the generic stabilizer is trivial. For the general $x \in X_{r}$, therefore, $\sum_{g \in G_{x}} f_{1-k \nu}(g)=1$ irrespective of $k \boldsymbol{\nu}$. If $k \boldsymbol{\nu}$ is even (respectively, odd) then there is a copy of $V_{k \boldsymbol{\nu}}$ in $H_{l}\left(X_{r}\right)$ if and only if $l$ is odd (respectively, even) and satisfies (2.1). Thus the number of copies of $V_{k \nu}$ in $H\left(X_{r}\right)$ is $\sim k \nu /\left(r^{2}-1\right)$, so that the dimension of $H_{k \nu}(\mathrm{X})$ is $\sim(k \nu)^{2}=\left(r^{2}-1\right)$.

### 2.2 Example about $G=U(2)$

Let $A$ be the hyperplane line bundle on $M=\mathbb{P}^{3}$; then the unit circle bundle $X \subseteq A^{\vee} \backslash\{\mathbf{0}\}$ may be identified with $S^{7} \subset \mathbb{C}^{4} \backslash\{\mathbf{0}\}$, and the projection $\pi: X \rightarrow \mathbb{P}^{3}$ with the Hopf map.

Consider the unitary representation of $U(2)$ on $\mathbb{C}^{4} \equiv \mathbb{C}^{2} \oplus \mathbb{C}^{2}$ given by

$$
\begin{equation*}
A \cdot(Z, W)=(A Z, A W) \tag{2.2}
\end{equation*}
$$

here $Z=\left(z_{1}, z_{2}\right)^{t}$, $W=\left(w_{1}, w_{2}\right)^{t} \in \mathbb{C}^{2}$. This linear action yields by restriction a contact action $\mu: G \times S^{7} \rightarrow S^{7}$, and descends to an holomorphic action $\mu: G \times \mathbb{P}^{3} \rightarrow \mathbb{P}^{3}$. If $\omega_{F S}$ is the Fubini-Study form on $\mathbb{P}^{3}$, then $\mu$ is Hamiltonian with respect to $2 \omega_{F S}$. The moment map is

$$
\Phi_{G}:[Z: W] \in \mathbb{P}^{3} \mapsto \frac{i}{\|Z\|^{2}+\|W\|^{2}}\left[z_{i} \bar{z}_{j}+w_{i} \bar{w}_{j}\right] \in \mathfrak{g} .
$$

Furthermore, $\tilde{\mu}$ is the contact lift of $\mu$.
From this, one can draw the following conclusions:
Lemma 2.2.1. Under the previous assumptions, we have:

1. $-i \Phi_{G}([Z: W])$ is a convex linear combination of the orthogonal projections onto the subspaces of $\mathbb{C}^{2}$ spanned by $Z$ and $W$, respectively;
2. $-i \Phi_{G}([Z: W])$ has rank 2 if and only if $Z$ and $W$ are linearly independent, rank 1 otherwise;
3. $\Phi_{G}(M)=i K$, where $K$ denotes the set of all positive semi-definite Hermitian matrices of trace 1;
4. the determinant of $-i \Phi_{G}([Z: W])$ is

$$
\operatorname{det}\left(-i \Phi_{G}([Z: W])\right)=\frac{|Z \wedge W|^{2}}{\left(\|Z\|^{2}+\|W\|^{2}\right)^{2}},
$$

where $Z \wedge W=z_{1} w_{2}-z_{2} w_{1} \in \mathbb{C}$;
5. the eigenvalues of $-i \Phi_{G}([Z: W])$ are both real and given by

$$
\lambda_{1,2}([Z: W])=\frac{1}{2}\left(1 \pm \sqrt{1-\frac{4|Z \wedge W|^{2}}{\left(\|Z\|^{2}+\|W\|^{2}\right)^{2}}}\right)
$$

Let us fix $\boldsymbol{\nu} \in \mathbb{Z}^{2}$ with $\nu_{1}>\nu_{2} \geq 0$. Let, as above, $\mathcal{O}_{\nu} \subseteq \mathfrak{g}$ denote the coadjoint orbit od $i D_{\nu}$. With $M=\mathbb{P}^{3}$, the locus $M_{\mathcal{O}_{\nu}}^{G}=\Phi_{G}^{-1}\left(\mathbb{R}_{+} \cdot \mathcal{O}_{\nu}\right)$ is given by the condition

$$
\nu_{2} \lambda_{1}([Z: W])-\nu_{1} \lambda_{2}([Z: W])=0 .
$$

In view of Lemma 2.2.1, this implies:
Corollary 2.2.1. Under the previous hypothesis,

$$
M_{\mathcal{O}_{\nu}}^{G}=\left\{[Z: W] \in \mathbb{P}^{3}: \frac{|Z \wedge W|}{\|Z\|^{2}+\|W\|^{2}}=\frac{\sqrt{\nu_{1} \nu_{2}}}{\nu_{1}+\nu_{2}}\right\}
$$

Let us now consider transversality. By Lemma 4.5.1 below (see also the discussion in $\S 2$ of [Pao12]), $\Phi_{G}$ is transverse to the ray $\mathbb{R}_{+} \cdot i D_{\nu}$ in $\mathfrak{g}$ if and only if $\tilde{\mu}$ is locally free along $X_{\nu}^{G}$ (that is, each $x \in X_{\nu}^{G}$ has discrete stabilizer).

On the other hand, by (2.2) $\tilde{\mu}$ is locally free at $(Z, W) \in S^{7}$ if and only if $Z \wedge W \neq 0$, and this is equivalent to $\Phi_{G}([Z: W])$ having rank 2; this means that $-i \Phi_{G}([Z: W])$ has two positive eigenvalues. Thus we obtain the following.

Corollary 2.2.2. The following conditions are equivalent:

1. $\Phi_{G}$ is transverse to $\mathbb{R}_{+} \cdot i D_{\nu}$ and $\Phi_{G}^{-1}\left(\mathbb{R}_{+} \cdot i D_{\nu}\right) \neq \emptyset$;
2. $\Phi_{G}$ is transverse to $\mathcal{O}_{\nu}$, and $\Phi_{G}^{-1}\left(\mathbb{R}_{+} \cdot \mathcal{O}_{\nu}\right) \neq \emptyset$;
3. $\nu_{1}, \nu_{2}>0$.

Let us now consider the restricted Hamiltonian action of $T$. Identifying $\mathfrak{t}$ with $i \mathbb{R}^{2}, \Phi_{T}: M \rightarrow \mathfrak{t}$ may be written:

$$
\begin{equation*}
\Phi_{T}: \mathbb{P}^{3} \ni[Z: W] \mapsto \frac{i}{\|Z\|^{2}+\|W\|^{2}}\binom{\left|z_{1}\right|^{2}+\left|w_{1}\right|^{2}}{\left|z_{2}\right|^{2}+\left|w_{2}\right|^{2}} \in \mathfrak{t} . \tag{2.3}
\end{equation*}
$$

Thus we obtain
Lemma 2.2.2. Assume that $\nu_{1}>\nu_{2} \geq 0$; then:

1. the image of $\Phi_{T}$ in $\mathfrak{t} \cong i \mathbb{R}^{2}$ is

$$
\Phi_{T}(M)=i\left\{\binom{x}{y}: x+y=1, x, y \geq 0\right\}
$$

2. the locus $M_{\nu}^{T}$ is given by

$$
M_{\nu}^{T}=\left\{[Z: W] \in \mathbb{P}^{3}: \nu_{2}\left(\left|z_{1}\right|^{2}+\left|w_{1}\right|^{2}\right)=\nu_{1}\left(\left|z_{2}\right|^{2}+\left|w_{2}\right|^{2}\right)\right\}
$$

3. $\Phi_{T}$ is transverse to $\mathbb{R}_{+} \cdot i D_{\nu}$ and $M_{\nu}^{G} \neq \emptyset$ if and only if $\nu_{1}, \nu_{2}>0$.

Proof. The first two statements follow immediately from (2.3). As to the third, let us recall again that $\Phi_{T}$ is transverse to $\mathbb{R}_{+} \cdot i D_{\nu}$ if and only if the action of $T$ on $X_{\nu}^{T} \subseteq S^{7}$ is locally free, [Pao12]. On the other hand, $T$ acts locally freely at $(Z, W) \in S^{7}$ if and only if $Z$ and $W$ are neither both scalar multiples of $\mathbf{e}_{1}$, nor both scalar multiples of $\mathbf{e}_{2}$, where $\left(\mathbf{e}_{1}, \mathbf{e}_{2}\right)$ is the standard basis of $\mathbb{C}^{2}$. By 2 ), there are no points $(Z, W)$ of this form in $X_{\nu}^{T}$ if and only if $\nu_{2}>0$.

Hence if $\nu_{1}, \nu_{2}>0$, then $\Phi_{G}$ and $\Phi_{T}$ are transverse to $\mathbb{R}_{+} \cdot \boldsymbol{\nu}$, and $M_{\nu}^{G} \neq \emptyset$, $M_{\nu}^{T} \neq \emptyset$. For instance,

$$
\left[\sqrt{\frac{\nu_{1}}{\nu_{1}+\nu_{2}}} \boldsymbol{e}_{1}: \sqrt{\frac{\nu_{1}}{\nu_{1}+\nu_{2}}} \boldsymbol{e}_{2}\right] \in M_{\nu}^{G} \cap M_{\nu}^{T} .
$$

More generally, we have the following.
Lemma 2.2.3. For any $\boldsymbol{\nu}, M_{\nu}^{G} \cap M_{\nu}^{T}=\Phi_{G}^{-1}\left\{i\left(\nu_{1}+\nu_{2}\right)^{-1} D_{\nu}\right\}$.
Proof. By Lemma 2.2.1, $[Z: W] \in M_{\nu}^{G}$ if and only if $-i \Phi_{G}([Z: W])$ is similar to $D_{\boldsymbol{\nu} /\left(\nu_{1}+\nu_{2}\right)}$; on the other hand, by Lemma 2.2.2, $[Z: W] \in M_{\boldsymbol{\nu}}^{G}$ if and only if for some $z \in \mathbb{C}$

$$
-i \Phi_{G}([Z: W])=\left(\begin{array}{cc}
\nu_{1} /\left(\nu_{1}+\nu_{2}\right) & z \\
\bar{z} & \nu_{1} /\left(\nu_{1}+\nu_{2}\right)
\end{array}\right) .
$$

Equating determinants, we conclude that $z=0$. This concludes the proof.

Let $\mathfrak{g}_{i} \subset \mathfrak{g}$ be the affine hyperplane of the skew-Hermitian matrices of trace $i$; we may interpret $\Phi_{G}$ as a smooth map $\Phi_{G}^{\prime}: \mathbb{P}^{3} \rightarrow \mathfrak{g}_{i}$.
Lemma 2.2.4. If $\nu_{1}>\nu_{2}>0$, then $i\left(\nu_{1}+\nu_{2}\right)^{-1} D_{\nu} \in \mathfrak{g}_{i}$ is a regular value of $\Phi_{G}^{\prime}$.
Proof. Clearly, the latter matrix is a regular value of $\Phi_{G}^{\prime}$ if and only if $\Phi_{G}$ is transverse to the ray $\mathbb{R}_{+} \cdot i D_{\nu}$; thus the statement follows from Corollary 2.2.2.

By Lemmata 2.2.3 and 2.2.4, we obtain
Corollary 2.2.3. Suppose $\nu_{1}>\nu_{2}>0$. Then, with $M=\mathbb{P}^{3}$ :

1. $M_{\mathcal{O}_{\nu}}^{G}$ and $M_{\nu}^{T}$ are smooth compact (real) hypersurfaces in $M$;
2. $\quad M_{\mathcal{O}_{\nu}}^{G} \cap M_{\nu}^{T}$ is a smooth submanifold of $M$ of real codimension 3.

Let us now describe the saturation $G \cdot M_{\nu}^{T}$.
Lemma 2.2.5. Under the previous assumptions

$$
G \cdot M_{\nu}^{T}=\left\{[Z: W] \in \mathbb{P}^{3}: \frac{\|Z \wedge W\|}{\|Z\|^{2}+\|W\|^{2}} \leq \frac{\sqrt{\nu_{1} \nu_{2}}}{\nu_{1}+\nu_{2}}\right\}
$$

Proof. Consider $[Z: W] \in \mathbb{P}^{3}$ with $(Z, W) \in S^{7}$. By definition, $[Z: W] \in$ $G \cdot M_{\nu}^{T}$ if and only if there exists $A \in G$ such that $[A Z: A W] \in M_{\nu}^{T}$; we may actually require without loss that $A \in S U(2)$. Let us write

$$
A=\left(\begin{array}{cc}
a & -\bar{c} \\
c & \bar{a}
\end{array}\right) \in S U(2), \quad Z=\binom{z_{1}}{z_{2}}, \quad W=\binom{w_{1}}{w_{2}} ;
$$

then $[A Z: A W] \in M_{\nu}^{T}$ if and only if (with some computations)

$$
\begin{aligned}
0 & =\nu_{2}\left(\left|a z_{1}-\bar{c} z_{2}\right|^{2}+\left|a w_{1}-\bar{c} w_{2}\right|^{2}\right)-\nu_{1}\left(\left|c z_{1}+\bar{a} z_{2}\right|^{2}+\left|c w_{1}+\bar{a} w_{2}\right|^{2}\right) \\
& =\nu_{2}\left\|\left(\begin{array}{cc}
z_{1} & z_{2} \\
w_{1} & w_{2}
\end{array}\right)\binom{a}{-\bar{c}}\right\|^{2}-\nu_{1}\left\|\left(\begin{array}{cc}
z_{1} & z_{2} \\
w_{1} & w_{2}
\end{array}\right)\binom{c}{\bar{a}}\right\|^{2}
\end{aligned}
$$

In other words, $[Z: W] \in G \cdot M_{\nu}^{T}$ if and only if there exits an orthonormal basis $\mathcal{B}=\left(V_{1}, V_{2}\right)$ of $\mathbb{C}^{2}$ such that

$$
\nu_{2}\left\|\left(\begin{array}{cc}
z_{1} & z_{2}  \tag{2.4}\\
w_{1} & w_{2}
\end{array}\right) V_{1}\right\|^{2}=\nu_{1}\left\|\left(\begin{array}{cc}
z_{1} & z_{2} \\
w_{1} & w_{2}
\end{array}\right) V_{2}\right\|^{2}
$$

Now, for any $V \in S^{7}$ we have

$$
\left\|\left(\begin{array}{cc}
z_{1} & z_{2} \\
w_{1} & w_{2}
\end{array}\right) V\right\|^{2}=V^{t}\left(\begin{array}{cc}
z_{1} & w_{1} \\
z_{2} & w_{2}
\end{array}\right)\left(\begin{array}{cc}
\bar{z}_{1} & \bar{z}_{2} \\
\bar{w}_{1} & \bar{w}_{2}
\end{array}\right) \bar{V}=V^{t} \frac{1}{i} \Phi_{G}([Z: W]) \bar{V} .
$$

If $\lambda_{1}(Z, W) \geq \lambda_{2}(Z, W) \geq 0$ are the eigenvalues of $-i \Phi_{G}([Z: W])$ (see Lemma 2.2.1), we then obtain for any $V \in S^{7}$

$$
\lambda_{1}(Z, W) \geq\left\|\left(\begin{array}{cc}
z_{1} & z_{2}  \tag{2.5}\\
w_{1} & w_{2}
\end{array}\right) V\right\|^{2} \geq \lambda_{2}(Z, W)
$$

with left (respectively, right) equality holding if and only if $V$ is an eigenvector of $-i \Phi_{G}([Z: W])$ relative to $\lambda_{1}(Z, W)$ (respectively $\lambda_{2}(Z, W)$ ). We conclude from (2.4) and (2.5) that if $(Z, W) \in G \cdot X_{\nu}^{T}$ then the following inequalities holds:

$$
\begin{equation*}
\nu_{1} \lambda_{1}(Z, W) \geq \nu_{2} \lambda_{2}(Z, W), \quad \nu_{2} \lambda_{1}(Z, W) \geq \nu_{1} \lambda_{2}(Z, W) \tag{2.6}
\end{equation*}
$$

While the former is trivial, since $\nu_{1}>\nu_{2}>0$ and $\lambda_{1}(Z, W) \geq \lambda_{2}(Z, W) \geq 0$, the latter is equivalent to the other

$$
\begin{equation*}
\frac{\sqrt{\nu_{1} \nu_{2}}}{\nu_{1}+\nu_{2}} \geq\|Z \wedge W\| \tag{2.7}
\end{equation*}
$$

Suppose, conversely, that (2.7) holds. Then (2.6) also holds. Let ( $W_{1}, W_{2}$ ) be an orthonormal basis of eigenvectors of $-i \Phi_{G}([Z: W])$ with respect to the eigenvalues $\lambda_{1}(Z, W)$ and $\lambda_{2}(Z, W)$, respectively. Evaluating the two sides of (2.4) with $V_{1}^{\prime}=W_{1}$ and $V_{2}^{\prime}=W_{2}$ in place of $\left(V_{1}, V_{2}\right)$. We obtain

$$
\nu_{2}\left\|\left(\begin{array}{cc}
z_{1} & z_{2} \\
w_{1} & w_{2}
\end{array}\right) V_{1}^{\prime}\right\|^{2}=\nu_{2} \lambda_{1}(Z, W) \geq \nu_{1} \lambda_{2}(Z, W)=\nu_{1}\left\|\left(\begin{array}{cc}
z_{1} & z_{2} \\
w_{1} & w_{2}
\end{array}\right) V_{2}^{\prime \prime}\right\|^{2}
$$

Using instead $V_{1}^{\prime \prime}=W_{2}$ and $V_{2}^{\prime \prime}=W_{1}$ in place of $\left(V_{1}, V_{2}\right)$, we obtain

$$
\nu_{2}\left\|\left(\begin{array}{cc}
z_{1} & z_{2} \\
w_{1} & w_{2}
\end{array}\right) V_{1}^{\prime \prime}\right\|^{2}=\nu_{2} \lambda_{2}(Z, W) \leq \nu_{1} \lambda_{1}(Z, W)=\nu_{1}\left\|\left(\begin{array}{cc}
z_{1} & z_{2} \\
w_{1} & w_{2}
\end{array}\right) V_{2}^{\prime \prime}\right\|^{2}
$$

Since $G=U(2)$ is connected, and acts transitively on the family of all orthonormal basis of $\mathbb{C}^{2}$, we conclude by continuity that there exists an orthonormal basis ( $V_{1}, V_{2}$ ) on which (2.4) is satisfied.

In view of Corollary 2.2.1, we deduce
Corollary 2.2.4. $M_{\mathcal{O}_{\nu}}^{G}=\partial\left(G \cdot M_{\nu}^{T}\right)$.
The boundary $\partial\left(G \cdot M_{\nu}^{T}\right)$ consists of those $[Z: W] \in \mathbb{P}^{3}$ such that $-i \Phi_{G}([Z: W])$ is similar to $\left(\nu_{1}+\nu_{2}\right)^{-1} D_{\nu}$, while the interior $\left(G \cdot M_{\nu}^{T}\right)^{0}$ consists of those $[Z: W] \in \mathbb{P}^{3}$ such that $-i \Phi_{G}([Z: W])$ is similar to a matrix of the form

$$
\frac{1}{\nu_{1}+\nu_{2}}\left(\begin{array}{cc}
\nu_{1} & z \\
\bar{z} & \nu_{2}
\end{array}\right),
$$

for some complex number $z \neq 0$.

Finally, the locus $X^{\prime} \subseteq X=S^{7}$ of those $(Z, W)$ at which $\tilde{\mu}$ is not locally free is defined by the condition $Z \wedge W=0$, and therefore it is contained in $\left(G \cdot M_{\nu}^{T}\right)^{0}$. It is the unit circle bundle over a non-singular quadric hypersurface in $\mathbb{P}^{3}$. The stabilizer subgroup of $(Z, W) \in S^{7}$ is trivial if $Z \wedge W \neq 0$, and it is isomorphic to $S^{1}$ otherwise.

For any fixed $\boldsymbol{\nu}=\left(\nu_{1}, \nu_{2}\right) \in \mathbb{Z}^{2}$ with $\nu_{1}>\nu_{2}$, let consider how $V_{k \nu}$ appears in the isotypical decomposition of $H(X)$ under $\hat{\mu}$. The Hopf map $\pi: X=S^{7} \rightarrow \mathbb{P}^{3}$ is the quotient map for the standard action $r: S^{1} \times$ $S^{7} \rightarrow S^{7} \subset \mathbb{C}^{4}$, given by complex scalar multiplication. The corresponding unitary representation of $S^{1}$ on $H(X)$ yields an isotypical decomposition $H(X)=\bigoplus_{l \in \mathbb{Z}} H_{l}(X)$, where for $l \in \mathbb{N}$ we set

$$
H_{l}(X):=\left\{f \in H(X): f\left(e^{i \theta} x\right)=e^{i l \theta} f(x) \forall x=(Z, W) \in X, e^{i \theta} \in S^{1}\right\}
$$

As is well-known, there are natural $U(2)$-equivariant unitary isomorphisms

$$
\begin{aligned}
H_{l}(X) & \cong H^{0}\left(\mathbb{P}^{3}, \mathcal{O}_{\mathbb{P}^{3}}(l)\right) \cong \operatorname{Sym}^{l}\left(\mathbb{C}^{2} \oplus \mathbb{C}^{2}\right) \\
& =\bigoplus_{h=0}^{l} \operatorname{Sym}^{h}\left(\mathbb{C}^{2}\right) \otimes \operatorname{Sym}^{l-h}\left(\mathbb{C}^{2}\right)
\end{aligned}
$$

On the other hand, a character computation yields the following.
Lemma 2.2.6. For $p \geq q$.

$$
\operatorname{Sym}^{p}\left(\mathbb{C}^{2}\right) \otimes \operatorname{Sym}^{q}\left(\mathbb{C}^{2}\right) \cong \bigoplus_{a=0}^{q}(\operatorname{det})^{\otimes a} \otimes \operatorname{Sym}^{p+q-2 a}\left(\mathbb{C}^{2}\right)
$$

as $U(2)$-representations.
Proof of Lemma 2.2.6. The character of $\operatorname{Sym}^{p}\left(\mathbb{C}^{2}\right)$ is $\chi_{(p+1,0)}$. Since the character of a tensor product of representations is the product of the respective characters, the character of $\operatorname{Sym}^{p}\left(\mathbb{C}^{2}\right) \otimes \operatorname{Sym}^{q}\left(\mathbb{C}^{2}\right)$ is $\chi^{\prime}:=\chi_{(p+1,0)} \cdot \chi_{(q+1,0)}$. Let us evaluate $\chi$ on a diagonal matrix $D_{\boldsymbol{z}}$ with diagonal $\boldsymbol{z}=\left(z_{1}, z_{2}\right)$. We
obtain

$$
\begin{align*}
\chi^{\prime}\left(D_{z}\right) & =\frac{z_{1}^{p+1}-z_{2}^{p+1}}{z_{1}-z_{2}} \cdot\left(z_{1}^{q}+z_{1}^{q-1} z_{2}+\cdots+z_{1} z_{2}^{q-1}+z_{2}^{q}\right)  \tag{2.8}\\
& =\frac{1}{z_{1}-z_{2}} \cdot\left(\sum_{j=0}^{q} z_{1}^{p+1+q-j} z_{2}^{j}-\sum_{j=0}^{q} z_{1}^{j} z_{2}^{p+1+q-j}\right) \\
& =\sum_{j=0}^{q} \frac{1}{z_{1}-z_{2}} \cdot\left(z_{1}^{p+1+q-j} z_{2}^{j}-z_{1}^{j} z_{2}^{p+1+q-j}\right) \\
& =\sum_{j=0}^{q} \chi_{p+1+q-j, j}\left(D_{\boldsymbol{z}}\right) .
\end{align*}
$$

Now, a character is uniquely determined by its restriction to $T$, and on the other hand the character of a direct sum is the sum of the characters; therefore, in view of (3.4), we conclude from (2.8) that

$$
\operatorname{Sym}^{p}\left(\mathbb{C}^{2}\right) \otimes \operatorname{Sym}^{q}\left(\mathbb{C}^{2}\right) \cong \bigoplus_{j=0}^{q} V_{(p+1+q-j, j)}=\bigoplus_{j=0}^{q} \operatorname{det}^{\otimes j} \operatorname{Sym}^{p+q-2 j}\left(\mathbb{C}^{2}\right)
$$

Therefore

$$
\begin{equation*}
H_{l}(X) \cong \bigoplus_{h=0}^{l} H_{l, h}(X) \tag{2.9}
\end{equation*}
$$

where we set

$$
\begin{equation*}
H_{l, h}(X):=\bigoplus_{a=0}^{\min (h, l-h)}(\operatorname{det})^{\otimes a} \otimes \operatorname{Sym}^{l-2 a}\left(\mathbb{C}^{2}\right) \tag{2.10}
\end{equation*}
$$

In order for the $a$-th summand in (2.9) to be isomorphic to $V_{k \nu}$, we need to have $a=k \nu_{2}$ and $l-2 a=k\left(\nu_{1}-\nu_{2}\right)-1$; hence in this special case $H(X)_{k \nu} \subseteq H_{l}(X)$ with $l=k\left(\nu_{1}+\nu_{2}\right)-1$. Let us estimate the multiplicity of $H(X)_{k \nu}$ in $H_{l}(X)$. In order for the $a$-th summand with $a=k \nu_{2}$ to appear in $H_{l h}(X)$ in (2.10) for some $h \leq k\left(\nu_{1}+\nu_{2}\right)-1$ we need to have

$$
\begin{aligned}
a=k \nu_{2} & \leq \min \left(h, k\left(\nu_{1}+\nu_{2}\right)-1-h\right) \\
& \Rightarrow k \nu_{2} \leq h, \quad k \nu_{2} \leq k\left(\nu_{1}+\nu_{2}\right)-1-h \\
& \Rightarrow k \nu_{2} \leq h \leq k \nu_{1}-1
\end{aligned}
$$

Hence there are $k\left(\nu_{1}-\nu_{2}\right)-1$ values of $h$ for which $H_{l, h}(X)$ contains one copy of $V_{k \nu}$. The dimension of $H(X)_{k \nu}$ is thus

$$
\left(k\left(\nu_{1}-\nu_{2}\right)-1\right) k\left(\nu_{1}-\nu_{2}\right) \sim k^{2}\left(\nu_{1}-\nu_{2}\right)^{2}+O(k) .
$$

## Chapter 3

## Preliminaries

### 3.1 Quantized manifolds

Recall that a $d$-dimensional complex projective manifold is a complex submanifolds of the complex projective space $\mathbb{P}^{n}(\mathbb{C})$. A positive complex line bundle on $M$ is a triple $(A, h, \nabla)$ where

1. $A$ is a complex line bundle,
2. $h$ is an Hermitian metric on $L$ (that is a smooth field of Hermitian inner products $\langle, \cdot\rangle_{h}$ in the fibers of $A$ );
3. $\nabla$ is a connection on $M$ compatible with $h$, that is

$$
\mathrm{d}\left\langle s_{1}, s_{2}\right\rangle_{h}=\left\langle\nabla s_{1}, s_{2}\right\rangle_{h}+\left\langle s_{1}, \nabla s_{2}\right\rangle_{h}
$$

for each smooth section $s_{1}, s_{2}$, such that the curvature of $\nabla$ (up to a factor $2 / i$ ) is a positive form $\omega$.
As a consequence $M$ is a Kähler manifold with Kähler form $\omega$.
From the bundle $(A, h)$ we pass to its dual $\left(A^{\vee}, h^{\vee}\right)$ and we consider the circle bundle

$$
X=\left\{a \in A: h^{\vee}(a, a)=1\right\}
$$

which is the boundary of the following strictly pseudoconvex domain (by the positivity of $(A, h))$

$$
D=\left\{a \in A: h^{\vee}(a, a) \leq 1\right\}
$$

Furthermore $X$ is a contact manifold, that is there exists a 1-form $\alpha$ on $X$ such that $\alpha \wedge(\mathrm{d} \alpha)^{n} \neq 0$ (see Example 3.5.11, p. 130 in [MS17]).

### 3.2 Compact Lie groups and actions

A Lie group $\left(G, \cdot{ }_{G}\right)$ is a group and a manifold such that the following maps ${ }^{\cdot}{ }_{G}: G \times G \rightarrow G$ and ${ }^{-1}: G \rightarrow G$ are smooths. Let us review some basics facts concerning Lie groups and its actions on symplectic manifolds. (We will write $g h$ instead of $g \cdot{ }_{G} h$ ).

### 3.2.1 Volume elements

We will denote by $\mathrm{dV}_{G}(g)$ the Haar measure on a compact Lie group $G$, which is the unique measure on $G$ invariant under both left and right translations, normalized by the condition

$$
\int_{G} \mathrm{dV}_{G}(g)=1
$$

For later use, we compute explicitly the volume element for $S U(2) / T$. There exists a diffeomorphism $\gamma$ from $S U(2)$ into the unit sphere $S^{3} \subset \mathbb{C}^{2}$,

$$
g:=\left(\begin{array}{cc}
\alpha & -\bar{\beta}  \tag{3.1}\\
\beta & \bar{\alpha}
\end{array}\right) \xrightarrow{\gamma}\binom{\alpha}{\beta} \in S^{3},
$$

where $\alpha$ and $\beta$ are complex numbers such that $|\alpha|^{2}+|\beta|^{2}=1$. Furthermore, $\gamma$ intertwines the right action of $T \cong S^{1}$ on $S U(2)$ with the standard circle action on $S^{3}$. Therefore, the projection $S U(2) \rightarrow S U(2) / T$ may be identified with the Hopf map $S^{3} \rightarrow \mathbb{P}^{1} \cong S^{2}$. It follows that the Haar measure on $S U(2) / T$ is a positive multiple of the pull-back of the standard measure on $S^{2}$.

In particular, when $\alpha \cdot \beta \neq 0$, we can set the following local coordinates on $S U(2)$ :

$$
\begin{equation*}
\alpha=e^{i \theta_{\alpha}} \cos \theta, \quad \beta=e^{i \theta_{\beta}} \sin \theta \tag{3.2}
\end{equation*}
$$

Thus, the Haar volume element on $S U(2) / T$ is

$$
\omega_{F S}=\frac{1}{2 \pi} \sin (2 \theta) d \theta d \delta
$$

### 3.2.2 Harmonic analysis on compact Lie groups

In this section we will resume some basic theorems in harmonic analysis of compact Lie groups $G$, in particular for $G=U(2)$ and $G=S U(2)$. We refer to the books [Var89] and [BtD95].

The central notion is that of the character of a finite dimensional unitary representation. Let us denote with $\kappa$ the representation, its character $\chi_{\kappa}$ is the function on $G$ defined by

$$
\chi_{\kappa}(g)=\operatorname{trace}(\kappa(g)), \quad g \in G .
$$

Let us recall that two representations $\kappa: G \rightarrow G L(V)$ and $\kappa^{\prime}: G \rightarrow$ $G L\left(V^{\prime}\right)$, respectively on the vector spaces $V$ and $V^{\prime}$, are called equivalent if there is a linear isomorphism $T: V \rightarrow V^{\prime}$ such that $\kappa^{\prime}(g)=T \kappa(g) T^{-1}$ for all $g \in G$. The character of $\kappa$ depends only on the equivalence class $[\kappa]$. The set of equivalence classes of irreducible representations of $G$ is written $\hat{G}$ and is called unitary dual of $G$.

Furthermore the characters are class functions, i.e. are invariant under the action of the group by conjugation. We will denote by $L^{2}(G)^{c}$ the space of square integrable class functions.

In addition to the characters we can associate with a finite dimensional unitary representation $\kappa$ its matrix elements, namely, the functions

$$
\mathcal{M}_{a, b}^{(\kappa)}: g \rightarrow(\kappa(g) a, b) \quad(a, b \in V)
$$

Theorem 3.2.1 (Completeness Theorem of Peter-Weyl). The irreducible unitary representations of $G$ are all finite-dimensional and they separate the points of $G$. The irreducible characters form an orthonormal basis of $L^{2}(G)^{c}$, and $L^{2}(G)$ is the orthogonal direct sum of matrix elements.

We now want to show some explicit expressions for the characters of $U(2)$ and we will deduce from these the ones of $S U(2)$. One is able to write explicitly a formula for the restriction of $\chi_{\kappa}$ to a maximal compact abelian subgroup of $U(2)$, the torus $T$. Let $d V_{T}$ be the Haar measure on $T$ and define the function $\Delta: T \rightarrow \mathbb{C}$ by setting

$$
\Delta(t):=t_{1}-t_{2}, \quad t:=\left(t_{1}, t_{2}\right) \in T
$$

here we identify $T$ with $S^{1} \times S^{1}$ in a natural manner. Furthermore, for any $f \in \mathcal{C}^{\infty}(U(2))$ let us define $F_{f}: T \rightarrow \mathbb{C}$ by setting

$$
F_{f}(t):=\int_{U(2) / T} f\left(g t g^{-1}\right) \mathrm{dV}_{U(2) / T}(g T)
$$

If $f$ is a class function $F_{f}(t)=f(t)$ for any $t \in T$.

Theorem 3.2.2 (Weyl's integration formula). A Borel function $f$ on $U(2)$ lies in $L^{1}(U(2))$ if and only if $F_{f}$ lies in $L^{1}\left(T, \Delta \bar{\Delta} \mathrm{dV}_{T}\right)$; in this case

$$
\int_{U(2)} f \mathrm{dV}_{U(2)}=\frac{1}{2} \int_{T} F_{f} \Delta \bar{\Delta} \mathrm{dV}_{T} .
$$

In particular, an invariant function $f$ on $U(2)$ lies in $L^{1}(U(2))$ if and only if its restriction to $T$ lies in $L^{1}\left(T, \Delta \bar{\Delta} \mathrm{dV}_{T}\right)$, and then

$$
\int_{U(2)} f \mathrm{dV}_{U(2)}=\left.\frac{1}{2} \int_{T} f\right|_{T} \Delta \bar{\Delta} \mathrm{dV}_{T}
$$

We will use Weyl's integration formula in the proof of Theorems described in the introduction and it is also the main tool in the proof of the following theorem, see [Var89].

Theorem 3.2.3 (Weyl's character formula). The irreducible characters of $U(2)$ are in one-one correspondence with the decreasing couples of integers. For each $\boldsymbol{\nu}=\left(\nu_{1}, \nu_{2}\right)$, with $\nu_{1}>\nu_{2}$, the character $\chi_{\nu}$ is given on $T$ by

$$
\begin{equation*}
\chi_{\nu}(t):=\frac{t_{1}^{\nu_{1}} t_{2}^{\nu_{2}}-t_{1}^{\nu_{2}} t_{2}^{\nu_{1}}}{t_{1}-t_{2}} \tag{3.3}
\end{equation*}
$$

Namely, $\boldsymbol{\nu}$ corresponds to the irreducible representation

$$
\begin{equation*}
\operatorname{det}^{\nu_{2}} \otimes \operatorname{Sym}^{\nu_{1}-\nu_{2}-1}\left(\mathbb{C}^{2}\right) \tag{3.4}
\end{equation*}
$$

where $\mathrm{Sym}^{\nu_{1}-\nu_{2}-1}$ is the space of homogeneous polynomials in two complex variables of degree $\nu_{1}-\nu_{2}-1$. Thus, the representations of $U(2)$ labeled by $\boldsymbol{\nu}$ has dimension $k\left(\nu_{1}-\nu_{2}\right)$.

The above theorems for $U(2)$ remain valid for $S U(2)$ with the sole modification that the character is $\chi_{(\nu, 0)}, \nu>0$. Explicitly, the Weyl character formula for $S U(2)$ is

$$
\begin{equation*}
\chi_{k \nu}\left(e^{i \vartheta}\right)=\frac{e^{i k \nu \vartheta}-e^{-i k \nu \vartheta}}{e^{i \vartheta}-e^{-i \vartheta}}=e^{i(k \nu-1) \vartheta}+e^{i(k \nu-3) \vartheta}+\cdots+e^{-i(k \nu-1) \vartheta} . \tag{3.5}
\end{equation*}
$$

Similarly every finite dimensional irreducible representation of $S U(2)$ is isomorphic to $\operatorname{Sym}^{\nu-1}\left(\mathbb{C}^{2}\right)$, for some $\nu>0$. Hence, the representations of $S U(2)$ labeled by $\boldsymbol{\nu}=(\nu, 0)$ has dimension $k \nu$.

### 3.2.3 Hamiltonian group actions

Let $M$ be a manifold, $G$ a compact connected Lie group and $\mathfrak{g}$ its Lie algebra. Let us suppose that the smooth map $\mu: G \times M \rightarrow M$ defines an action of $G$ on $M$. For each point $m \in M$, the mapping

$$
\begin{aligned}
f_{m}: G & \rightarrow M \\
& g \mapsto \mu_{g}(m):=g \cdot m
\end{aligned}
$$

is a smooth map and its differential

$$
\begin{aligned}
\mathrm{d}_{m} f: \mathfrak{g} & \rightarrow T_{m} M \\
\xi & \mapsto \xi_{M}(m)
\end{aligned}
$$

is a Lie algebra homeomorphism. The vector field $\xi_{M}$ is called infinitesimal vector field associated to $\xi$. If we have defined on $M$ a complex line bundle on $A$ (with projection $p: A \rightarrow M$ ) then we say that $\mu$ can be lifted to a group action $\tilde{\mu}: G \times A \rightarrow A$ on $A$ if there exists a group action $\tilde{\mu}$ of $G$ on $A$ such that the induced restricted action on the fiber of $A$ is linear.

If $(M, \omega)$ is symplectic, we suppose that $G$ acts on $M$ via symplectomorphism $\mu_{g}: M \rightarrow M$, then for every $\xi \in \mathfrak{g}$ the infinitesimal vector field $\xi_{M}$ is symplectic. This means that $\iota\left(\xi_{M}\right) \omega$ is closed for every $\xi$; if the 1 -form $\iota\left(\xi_{M}\right) \omega$ is exact then we say that the action is weakly Hamiltonian. More explicitly there exists a function $\Phi^{\xi} \in \mathcal{C}^{\infty}(M)$ such that $\mathrm{d} \Phi^{\xi}=\iota\left(\xi_{M}\right) \omega$.

The group $G$ acts on itself by conjugation $C: G \times G \rightarrow G$. The differential of $C_{g}$ defines an action $\operatorname{Ad}_{g}: \mathfrak{g} \rightarrow \mathfrak{g}$ of $G$ on its Lie algebra. The dual $\mathrm{Ad}^{*}$ of Ad is called the co-adjoint action. The action is called Hamiltonian if the map $\xi \rightarrow \Phi^{\xi}$ can be chosen to be $G$-equivariant, which means

$$
\Phi^{\xi} \circ \mu_{g}=\Phi^{\operatorname{Ad}_{g} \xi}
$$

for every $\xi \in \mathfrak{g}$ and $g \in G$. This last condition is equivalent to require that the $\operatorname{map} \xi \rightarrow \Phi^{\xi}$ is a Lie algebra homomorphism with respect to the Poisson structure $\{$,$\} on \mathcal{C}^{\infty}(M)$ (see Lemma 5.2.1, p. 203 of [MS17]).

Now, assume that the action $\mu$ is Hamiltonian. Then a moment map $\Phi: M \rightarrow \mathfrak{g}^{\vee}$ is a $G$-equivariant smooth map such that $\Phi^{\xi}(m)=\langle\Phi(m), \xi\rangle$, where $\langle\cdot, \cdot\rangle$ is the pairing between $\mathfrak{g}$ and its dual. Equivalently,

$$
\mathrm{d} \Phi_{m}^{\xi}(\mathbf{v})=\omega_{m}\left(\mathbf{v}, \xi_{M}(m)\right)
$$

for each $m \in M, \mathbf{v} \in T_{m} M$ and $\xi \in \mathfrak{g}$.

### 3.3 The Szegö projector

Let $W \subseteq \mathbb{C}^{n}$ be a compact strictly pseudo-convex domain with smooth boundary. Let $r$ be a smooth function on $W$ with $r>0$ on $\operatorname{Int} W, r=0$ on $\partial W$ and $\mathrm{d} r \neq 0$ near $\partial W$. Let $\iota: \partial W \rightarrow W$ be the inclusion map; the one-form $\alpha:=\iota^{*} \operatorname{Im} \bar{\partial} r$ is a contact form on $\partial W$. Thus, $\nu:=\alpha \wedge(\mathrm{d} \alpha)^{n-1}$ defines a volume form on $\partial W$. Let $L^{2}(\partial W, \nu)$ be the space of square integrable function over $\partial W$ with respect to the measure $\nu$. The Hardy Space is the closure in $L^{2}(\partial W, \nu)$ of the space of smooth function on $\partial W$ which can be extended to holomorphic function on $W$. The orthogonal projector $\Pi: L^{2}(\partial W, \nu) \rightarrow H^{2}(\partial W)$ is called the Szegö projector.

Let $M$ be a connected complex $d$-dimensional projective manifold with a $(A, h)$ a positive line bundle on $M$, as described in chapter 1.1. If $A^{\vee} \supset$ $X \xrightarrow{\pi} M$ is the unit circle bundle in the dual of $A$, then $X$ is the boundary of a strictly pseudo-convex domain and the structures defined in the previous paragraph carry over this case. Here the domain $W$ is the disc bundle,

$$
W:=\left\{(m, \nu) \in A^{\vee}: h(\nu, \nu) \leq 1\right\}
$$

and the defining function $r$ of the previous paragraph is $r(m, \nu)=1-\|\nu\|_{m}^{2}$, where $\|\cdot\|$ is the norm in the metric induced by $h$.

### 3.3.1 The Szegö kernel

Let us recall some basic facts concerning the Szegö kernel from [BdMS76], [Zel98]. The Szegö kernel is a Fourier integral operator of complex type, in particular there exists a symbol $s \in S^{d}\left(X \times X \times \mathbb{R}^{+}\right)$of the type

$$
s(x, y, u) \sim \sum_{k=0}^{+\infty} u^{d-k} s_{k}(x, y)
$$

to be so that

$$
\Pi(x, y)=\int_{0}^{+\infty} e^{i u \psi(x, y)} s(x, y, u) \mathrm{d} u
$$

where the phase $\psi \in \mathcal{C}^{\infty}(X \times X)$ is determined up to a function which vanishes to infinite order along the diagonal. Explicitly, $\psi$ is determined by the following properties:

1. $\psi(x, x)=\frac{1}{i} r(x)$, where $r$ is the defining function of $X$;
2. $\mathrm{d}_{x}^{\prime \prime} \psi$ and $\mathrm{d}_{y}^{\prime} \psi$ vanish to infinite order along the diagonal;
3. $\psi(x, y)=-\overline{\psi(y, x)}$.

Example 3.3.1. Consider the unit ball in $\mathbb{C}^{n+1}$. The above formula has the form

$$
\Pi(x, y)=\frac{1}{(1-\langle x, y\rangle)^{n+1}}=\int_{0}^{\infty} e^{i u \psi(x, y)} u^{n} \mathrm{~d} u
$$

where $\psi(x, y)=1-\langle x, y\rangle$.
In [SZ02], Shiffman and Zelditch provide some results about the scaling asymptotics of Szegö kernel and their implications in symplectic geometry. Their proofs rely on the notion of Heisenberg local coordinates. A set of Heisenberg local coordinates on $X$ centred at $x$ determines a linear isometry $T_{x} X \cong \mathbb{R} \oplus \mathbb{C}^{d}$, furthermore, they are horizontal at $x$ with respect to the connection 1-form.

### 3.3.2 Equivariant Szegö projector

Suppose we have an unitary action of a compact Lie group $G$ on the Hardy space $H(X)$. We have a natural decomposition, given by

$$
H(X)=\bigoplus_{\boldsymbol{\nu} \in \hat{G}} H_{\boldsymbol{\nu}}(X)
$$

Let us focus on the case $G=U(2)$. Let $\Pi_{\nu}$ the projector onto the isotypical component $H_{\nu}(X)$. We will denote its kernel by $\Pi_{\nu}(\cdot, \cdot)$. For any $x, y \in X$, we have an explicit formula for the equivariant projector (see [GS82c]):

$$
\begin{equation*}
\Pi_{\nu}(x, y)=d_{\boldsymbol{\nu}} \int_{G} \overline{\chi_{\nu}(g)} \Pi\left(\tilde{\mu}_{g^{-1}}(x), y\right) \mathrm{dV}_{G}(g) \tag{3.6}
\end{equation*}
$$

where $d_{\boldsymbol{\nu}}$ is the dimension of the representation labeled by $\boldsymbol{\nu}$. In order to prove the main results we need to write (3.6) more explicitly. In view of the Weyl's integration formula the expression (3.6) can be rewritten

$$
\begin{equation*}
\Pi_{\nu}(x, y)=\frac{d_{\nu}}{2} \int_{T} \overline{\chi_{\nu}(t)} \overline{\Delta(t)} \Delta(t) F(t ; x, y) \mathrm{dV}_{T}(t) \tag{3.7}
\end{equation*}
$$

where

$$
\begin{equation*}
F(t ; x, y):=\int_{G / T} \Pi\left(\tilde{\mu}_{g t^{-1} g^{-1}}(x), y\right) \mathrm{dV}_{G / T}(g T) \tag{3.8}
\end{equation*}
$$

Thus, inserting the Weyl's character formula (3.3) for $U(2)$ we obtain, instead of (3.7),

$$
\Pi_{\nu}(x, y)=\frac{d_{\nu}}{2} \int_{T}\left(t_{1}^{-\nu_{1}} t_{2}^{-\nu_{2}}-t_{1}^{-\nu_{2}} t_{2}^{-\nu_{1}}\right) \Delta(t) F(t ; x, y) \mathrm{dV}_{T}(t)
$$

The previous integral splits in two summands that, after a change of variables, add up to

$$
\begin{equation*}
\Pi_{\nu}(x, y)=d_{\nu} \int_{T} t^{-\boldsymbol{\nu}} \Delta(t) F(t ; x, y) \mathrm{dV}_{T}(t) \tag{3.9}
\end{equation*}
$$

where $t^{-\nu}=t_{1}^{-\nu_{1}} t^{-\nu_{2}}$ and $d_{\nu}$ is explicitly given by $d_{\nu}=\nu_{1}-\nu_{2}$. For $G=$ $S U(2)$ the computations are similar. In the proof of Theorems in the next chapters we will gloss over these computations starting directly with (3.9).

## Chapter 4

## Proofs

### 4.1 Proof of Theorem 1.2.1

Proof. We have

$$
\Pi_{k \nu}(x, y)=k \nu \int_{G} \overline{\chi_{k \nu}(g)} \Pi\left(\tilde{\mu}_{g^{-1}}(x), y\right) \operatorname{dV}_{G}(g) .
$$

Let $\delta>0$ and let $\rho: \mathbb{R} \rightarrow(0,+\infty)$ be a bump function with $\rho \equiv 1$ on $(-\delta, \delta)$ and $\rho \equiv 0$ on $(-2 \delta, 2 \delta)^{c}$. We have

$$
\Pi_{k \nu}(x, y)=\Pi_{k \nu}(x, y)^{\prime}+\Pi_{k \nu}(x, y)^{\prime \prime}
$$

where

$$
\begin{aligned}
& \Pi_{k \nu}(x, y)^{\prime} \\
& \quad=k \nu \int_{G} \rho\left(\operatorname{dist}_{X}\left(\tilde{\mu}_{g^{-1}}(x), y\right)\right) \overline{\chi_{k \nu}(g)} \Pi\left(\tilde{\mu}_{g^{-1}}(x), y\right) \mathrm{dV}_{G}(g),
\end{aligned}
$$

and

$$
\begin{aligned}
& \Pi_{k \nu}(x, y)^{\prime \prime} \\
& \quad=k \nu \int_{G}\left(1-\rho\left(\operatorname{dist}_{X}\left(\tilde{\mu}_{g^{-1}}(x), y\right)\right)\right) \overline{\chi_{k \nu}(g)} \Pi\left(\tilde{\mu}_{g^{-1}}(x), y\right) \mathrm{dV}_{G}(g) .
\end{aligned}
$$

Since the function

$$
g \mapsto\left(1-\rho\left(\operatorname{dist}_{X}\left(\tilde{\mu}_{g^{-1}}(x), y\right)\right)\right) \Pi\left(\tilde{\mu}_{g^{-1}}(x), y\right)
$$

is smooth on $G$, we obtain that $\Pi_{k \nu}(x, y)^{\prime \prime}=O\left(k^{-\infty}\right)$.
We need to exploit the explicit description of $\Pi$ as an FIO developed in [BdMS76]. Namely, up to a smoothing contribution, we have

$$
\begin{aligned}
& \Pi_{k \nu}(x, y) \\
& \quad \sim k \nu \int_{0}^{+\infty} \mathrm{d} u \int_{G} \mathrm{dV} V_{G}(g)\left[\rho_{x}(g) \overline{\chi_{k \nu}(g)} e^{i u \psi\left(\tilde{\mu}_{g^{-1}}(x), y\right)} s\left(\tilde{\mu}_{g^{-1}}(x), y, u\right)\right] \\
& \quad \sim k^{2} \nu \int_{0}^{+\infty} \mathrm{d} u \int_{G} \mathrm{dV} V_{G}(g)\left[\rho_{x}(g) \overline{\chi_{k \nu}(g)} e^{i k u \psi\left(\tilde{\mu}_{g^{-1}}(x), y\right)} s\left(\tilde{\mu}_{g^{-1}}(x), y, k u\right)\right]
\end{aligned}
$$

where we have set

$$
\rho_{x}(g):=\rho\left(\operatorname{dist}_{X}\left(\tilde{\mu}_{g^{-1}}(x), y\right)\right) .
$$

Now, let $D \gg 0$ and $\rho_{1}: \mathbb{R} \rightarrow[0,+\infty)$ such that $\rho_{1} \equiv 1$ on $[0, D]$ and $\rho \equiv 0$ on $[2 D,+\infty)$. We have

$$
\Pi_{k \nu}(x y) \sim \Pi_{k \nu}(x, y)_{1}+\Pi_{k \nu}(x, y)_{2}
$$

where

$$
\left.\begin{array}{rl}
\Pi_{k \nu}(x, y)_{1}:=k^{2} \nu & \int_{0}^{+\infty} \mathrm{d}
\end{array}\right) \int_{G} \mathrm{dV}_{G}(g), ~\left(\rho_{x}(g) \rho_{1}(u) \overline{\chi_{k \nu}(g)} e^{i k u \psi\left(\tilde{\mu}_{g^{-1}}(x), y\right)} s\left(\tilde{\mu}_{g^{-1}}(x), y, k u\right)\right],
$$

and

$$
\begin{aligned}
\Pi_{k \nu}(x, y)_{2}:= & k^{2} \nu \int_{0}^{+\infty} \mathrm{d} u \int_{G} \mathrm{dV} V_{G}(g) \\
& {\left[\rho_{x}(g)\left(1-\rho_{1}(u)\right) \overline{\chi_{k \nu}(g)} e^{i k u \psi\left(\tilde{\mu}_{g^{-1}}(x), y\right)} s\left(\tilde{\mu}_{g^{-1}}(x), y, k u\right)\right] . }
\end{aligned}
$$

Lemma 4.1.1. $\Pi_{k \nu}(x, y)_{2}=O\left(k^{-\infty}\right)$ as $k \rightarrow+\infty$.
Proof of Lemma 4.1.1. Let $\left\{G^{\prime}, G^{\prime \prime}\right\}$ be an open cover of $G$ where

$$
\begin{aligned}
G^{\prime} & :=\left\{g \in G: \operatorname{dist}_{G}\left(g,\left\{ \pm I_{2}\right\}\right)<2 \delta\right\}, \\
G^{\prime \prime} & :=\left\{g \in G: \operatorname{dist}_{G}\left(g,\left\{ \pm I_{2}\right\}\right)>\delta\right\} .
\end{aligned}
$$

We may consider a subordinate partition of unity $\beta^{\prime}+\beta^{\prime \prime}=1$ on $G$. Let us set

$$
\varrho^{\prime}:=\varrho \cdot \beta^{\prime}, \quad \varrho^{\prime \prime}:=\varrho \cdot \beta^{\prime \prime} .
$$

Then $\varrho=\varrho^{\prime}+\varrho^{\prime \prime}$, where $\varrho^{\prime}$ is supported in a small neighborhood of $\left\{ \pm I_{2}\right\}$, and $\varrho^{\prime \prime}$ is supported away from $\left\{ \pm I_{2}\right\}$.

Accordingly, we have

$$
\Pi_{k \nu}(x, y)_{2}=\Pi_{k \nu}(x, y)_{2}^{\prime}+\Pi_{k \nu}(x, y)_{2}^{\prime \prime}
$$

where in the former (respectively, latter) summand $\varrho(g)$ has been replaced by $\varrho^{\prime}(g)$ (respectively, $\varrho^{\prime \prime}(g)$ ). Let us study the two summands in the previous expression separately.
Lemma 4.1.2. $\Pi_{k \nu}(x, y)_{2}^{\prime}=O\left(k^{-\infty}\right)$ as $k \rightarrow+\infty$.
Proof of Lemma 4.1.2. Let us assume to fix ideas that $k \boldsymbol{\nu}=2 l+1$ is odd. Then $V_{k \nu}$ may be identified with the vector space $\mathbb{C}^{(2 l)}\left[z_{1}, z_{2}\right]$ of complex homogeneous polynomials of degree $2 l$ in two variables. A natural basis of the latter is given by the monomials $P_{\mu}\left(z_{1}, z_{2}\right):=z_{1}^{l-\mu} z_{2}^{l+\mu}$, where $\mu \in$ $\{-l, \ldots, 0, \ldots l\}$. We shall accordingly denote the matrix elements of the representation $V_{k \nu}$ by

$$
\mathcal{M}_{a, b}^{(k \nu)}: G \rightarrow \mathbb{C}, \quad(a, b \in\{-l, \ldots, 0, \ldots, l\})
$$

We have

$$
\chi_{k \nu}(g)=\sum_{a=-l}^{l} \mathcal{M}_{a, a}^{(k \nu)}(g) .
$$

Thus, we get

$$
\Pi_{k \nu}(x, y)_{2} \sim \sum_{a=-l}^{l} \Pi_{k \nu}(x, y)_{2, a}
$$

where

$$
\begin{aligned}
& \Pi_{k \nu}(x, y)_{2, a} \quad:=k^{2} \nu \int_{0}^{+\infty} \mathrm{d} u \int_{G} \mathrm{dV}_{G}(g) \\
& \quad\left[\rho_{x}(g)\left(1-\rho_{1}(u)\right) \frac{\mathcal{M}_{a, a}^{k \nu}(g)}{} e^{i k u \psi\left(\tilde{\mu}_{g^{-1}}(x), y\right)} s\left(\tilde{\mu}_{g^{-1}}(x), y, k u\right)\right] .
\end{aligned}
$$

On the support of $\varrho^{\prime}$, either $g \sim I_{2}$ or $g \sim-I_{2}$, and hence we can write explicatively

$$
g=\left(\begin{array}{cc}
A(g) e^{i \theta_{G}(g)} & -\overline{\gamma(g)}  \tag{4.1}\\
\gamma(g) & A(g) e^{-i \theta_{G}(g)}
\end{array}\right),
$$

where $A(g)>0$, and either $\theta_{G}(g) \approx 0$ or $\theta_{G}(g) \approx \pi$. Furthermore, $A, \gamma$ and $\theta_{G}$ are smooth functions of $g \in G$ on a neighborhood of the support of $\varrho^{\prime}$.

We can write (see [App14], §2.7, formula (2.6.40)), with $g$ as in (4.1), the following expression for the matrix elements

$$
\begin{equation*}
\mathcal{M}_{a, a}^{(k \nu)}(g)=e^{-2 i a \theta_{G}(g)} \cdot P_{l, a}\left(A(g)^{2}\right), \tag{4.2}
\end{equation*}
$$

where $P_{l, a}$ may be expressed in terms of suitable Jacobi polynomials, and is itself a real polynomial. Since the left hand side of (4.2) is an entry of a unitary matrix, we have at any rate $\left|P_{l, a}\left(A(g)^{2}\right)\right| \leq 1$.

Hence, we have

$$
\begin{gathered}
\Pi_{k \nu}(x, y)_{2, a}:=k^{2} \nu \int_{G} \mathrm{~d} V_{G}(g) \int_{D / 2}^{+\infty} \mathrm{d} u\left[e^{i k \Psi_{a / k}(x, y ; g, u)}\right. \\
\left.\varrho_{x}(g) \cdot \varrho_{2}^{\prime}(u) \cdot P_{l, a}\left(A(g)^{2}\right) s\left(\tilde{\mu}_{g^{-1}}(x), y, k u\right)\right]
\end{gathered}
$$

where we have set

$$
\Psi_{a / k}(x, y ; g, u):=u \cdot \psi\left(\tilde{\mu}_{g^{-1}}(x), y\right)+\frac{2 a}{k} \cdot \theta_{G}(g) .
$$

Since

$$
-\frac{\nu}{2} \leq \frac{a}{k} \leq \frac{\nu}{2},
$$

the family of phases $\Psi_{a / k}$ is finite.
Let us view $\beta$, the infinitesimal generator of the standard torus $T$, as a left-invariant vector field on $G$; the corresponding 1-parameter group of diffeomorphisms is $\varphi_{\tau}(g):=g e^{\tau \beta}$. Therefore, if $L_{\beta}$ is the same vector field viewed as a differential operator on $G$, then $L_{\beta}\left(\theta_{G}\right)=1$ on the support of $\varrho^{\prime}$. Also, $L_{\beta}$ is a skew-hermitian operator on $L^{2}(G)$, since $\varphi_{\tau}$ induces a 1-parameter group of unitary automorphisms of $L^{2}(G)$; hence, $L_{\beta}^{t}=-\bar{L}_{\beta}$. Furthermore, the function $g \mapsto A(g)^{2}$ is smooth, real and $\varphi_{\tau}$-invariant; therefore, $L_{\beta}\left(A(g)^{2}\right)=\bar{L}_{\beta}\left(A(g)^{2}\right)=0$.

On the other hand, we have

$$
\begin{aligned}
\left.\frac{\mathrm{d}}{\mathrm{~d} \tau} \tilde{\mu}_{\varphi_{\tau}(g)^{-1}}(x)\right|_{\tau=0} & =\left.\frac{\mathrm{d}}{\mathrm{~d} \tau} \tilde{\mu}_{e^{-\tau \beta}}\left(\tilde{\mu}_{g^{-1}}(x)\right)\right|_{\tau=0} \\
& =-\beta_{X}\left(\tilde{\mu}_{g^{-1}}(x)\right) \\
& =-\beta_{M}\left(\tilde{\mu}_{g^{-1}}\left(m_{x}\right)\right)^{\sharp}+\left\langle\Phi_{G}\left(\mu_{g^{-1}}\left(m_{x}\right)\right), \beta\right\rangle \partial_{\theta}
\end{aligned}
$$

For $\varrho(g) \neq 0$ we have $\operatorname{dist}_{X}\left(\tilde{\mu}_{g^{-1}}(x), y\right) \leq 2 \delta$; therefore,

$$
\begin{equation*}
\left\langle\Phi_{G}\left(\mu_{g^{-1}}\left(m_{x}\right)\right), \beta\right\rangle=\left\langle\Phi_{G}\left(m_{y}\right), \beta\right\rangle+O(\delta) . \tag{4.3}
\end{equation*}
$$

If $\delta$ is sufficiently small, (4.3) is non-zero, since we are assuming that $\Phi_{G}\left(m_{y}\right)$ is not anti-diagonal; assuming to fix ideas that $\Phi_{G}\left(m_{y}\right)$ is diagonal, then the right hand side of $(4.3)$ is $\left\langle\Phi_{G}\left(m_{y}\right), \beta\right\rangle=2 \lambda\left(m_{x}\right)+O(\delta)$.

In addition, by the discussion in [SZ02], where $\varrho(g) \neq 0$

$$
\mathrm{d}_{\left(\tilde{\mu}_{g^{-1}}(x), y\right)} \psi=\left(\alpha_{\tilde{\mu}_{g^{-1}}(x)},-\alpha_{y}\right)+O(\delta) .
$$

Therefore,

$$
L_{\beta}\left(\Psi_{a / k}(x ; u, g)\right)=2\left[u \cdot \lambda\left(m_{x}\right)+\frac{a}{k}\right]+O(\delta) .
$$

For $u \gg 0$, we conclude that $L_{\beta}\left(\Psi_{a / k}(u, g)\right) \geq C^{\prime} \cdot u+1$ for some $C^{\prime}>0$, which can be chosen uniformly for all $a \in\{-l, \ldots, 0, \ldots l\}$; by iteratively integrating by parts by the transpose operator $L_{\beta}^{t}=-\bar{L}_{\beta}$, we conclude that $\Pi_{k \nu}(x, y)_{2, a}^{\prime}$ is rapidly decreasing as $k \rightarrow+\infty$ uniformly for $a \in\{-l \ldots, l\}$. Since this holds uniformly for each of the $k \boldsymbol{\nu}$ summands, the statement of Lemma is established in the case where $k \boldsymbol{\nu}$ is odd.

The case where $k \boldsymbol{\nu}$ is even is only slightly different - one takes $l$ to be half-integer (see Theorem 11.7.1 of [RT10]).

Lemma 4.1.3. $\Pi_{k \nu}(x, y)_{2}^{\prime \prime}=O\left(k^{-\infty}\right)$ as $k \rightarrow+\infty$.
Proof of Lemma 4.1.3. The proof is similar to the one of previous Lemma, except that we shall use eigenvalues rather than matrix elements, so we'll be somewhat sketchy.

If $g \in G \backslash\left\{ \pm I_{2}\right\}$, then there is a unique $\vartheta_{G}(g)=\cos ^{-1}(\operatorname{trace}(g) / 2) \in(0, \pi)$ such that the eigenvalues of $g$ are $e^{ \pm i \vartheta_{G}(g)}$. The map $g \in G \backslash \pm I_{2} \mapsto \vartheta_{G}(g) \in$ $(0, \pi)$ is $\mathcal{C}^{\infty}$. On the same domain, the character of $V_{k \nu}$ is thus given by

$$
\chi_{k \nu}(g)=\sum_{j=0}^{k \nu-1} e^{(k \nu-1-2 j) \vartheta_{G}(g)} .
$$

We shall now write

$$
\Pi_{k \nu}(x, y)_{2}^{\prime \prime} \sim \sum_{j=0}^{k \nu-1} \Pi_{k \nu}(x, y)_{2, j}^{\prime \prime},
$$

where $\Pi_{k \nu}(x, y)_{2, j}^{\prime \prime}$ has the same integrand as $\Pi_{k \nu}(x, y)_{2}^{\prime \prime}$ with $\chi_{k \nu}(g)$ replaced by $e^{(k \nu-1-2 j) \vartheta_{G}(g)}$. Thus we can see $\Pi_{k \nu}(x, y)_{2, j}^{\prime \prime}$ as a sum of oscillatory integrals with phases

$$
\Psi_{b}(x, y ; u, g)=u \psi\left(\tilde{\mu}_{g^{-1}}(x), y\right)-(\nu+b) \cdot \vartheta_{G}(g),
$$

where $b \in \mathbb{R}$ and again the phases $\Psi_{b}$ vary in a bounded family. Notice that, away from $\left\{ \pm I_{2}\right\}, \vartheta_{G}: G \rightarrow \mathbb{C}$ is a smooth function. We can argue as in the proof of the previous Lemma.

Thus, we have

$$
\Pi_{k \nu}(x, y) \sim \Pi_{k \nu}(x, y)_{1}
$$

Let us shrink further the support of the integrand. Let $\rho_{2}: \mathbb{R} \rightarrow[0,+\infty)$ be a smooth function such that $\rho_{2} \equiv 0$ on $(0,1 /(2 D)]$ and $\rho_{2} \equiv 1$ on $[1 / D,+\infty)$. We have

$$
\Pi_{k \nu}(x, y) \sim \Pi_{k \nu}(x, y)_{11}+\Pi_{k \nu}(x, y)_{12}
$$

where

$$
\begin{aligned}
\Pi_{k \nu}(x, y)_{11}:= & k^{2} \nu \int_{0}^{+\infty} \mathrm{d} u \int_{G} \mathrm{~d} V_{G}(g) \\
& {\left[\rho_{x}(g) \rho_{2}(u) \rho_{1}(u) \overline{\chi_{k \nu}(g)} e^{i k u \psi\left(\tilde{\mu}_{g^{-1}}(x), y\right)} s\left(\tilde{\mu}_{g^{-1}}(x), y, k u\right)\right] }
\end{aligned}
$$

and

$$
\begin{aligned}
& \Pi_{k \nu}(x, y)_{12}:=k^{2} \nu \int_{0}^{+\infty} \mathrm{d} u \int_{G} \mathrm{dV}_{G}(g) \\
& \quad\left[\rho_{x}(g)\left(1-\rho_{2}(u)\right) \rho_{1}(u) \overline{\chi_{k \nu}(g)} e^{i k u \psi\left(\tilde{\mu}_{g^{-1}}(x), y\right)} s\left(\tilde{\mu}_{g^{-1}}(x), y, k u\right)\right] .
\end{aligned}
$$

Lemma 4.1.4. $\Pi_{k \nu}(x x)_{12}=O\left(k^{-\infty}\right)$ as $k \rightarrow+\infty$.
Proof. By the use of the Weyl integration formula we have

$$
\begin{aligned}
\Pi_{k \nu}(x, y)_{12}=\frac{k^{2} \nu}{2 \pi} \int_{0}^{+\infty} & \mathrm{d} u \int_{G / T} \mathrm{~d} V_{G / T}(g T) \int_{-\pi}^{\pi} \mathrm{d} \vartheta \\
& {\left[\rho_{x}\left(g e^{i \vartheta} g^{-1}\right)\left(1-\rho_{2}(u)\right) \rho_{1}(u)\left(e^{i \vartheta}-e^{-i \vartheta}\right) e^{-i k \nu \vartheta}\right.} \\
& \left.e^{i k u \psi\left(\tilde{\mu}_{g e^{-i} \vartheta_{g^{-1}}(x), x}\right.} s\left(\tilde{\mu}_{g e^{-i \vartheta} g_{g^{-1}}}(x), x, k u\right)\right] .
\end{aligned}
$$

Let us set

$$
\Psi(u, g T, \vartheta)=u \psi\left(\tilde{\mu}_{g e^{-i} \vartheta_{g^{-1}}}(x), y\right)-\nu \vartheta
$$

If $D \gg 0$ and $0<u<1 / D$ we have

$$
\left|\partial_{\vartheta} \Psi(u, g T, \vartheta)\right| \geq \frac{\nu}{2}
$$

The statement follows integrating by parts in $\mathrm{d} \vartheta$.
In view of the previous lemmas, we can write

$$
\begin{align*}
& \Pi_{k \nu}(x, y) \sim \frac{k^{2} \nu}{2 \pi} \int_{1 / D}^{D} \mathrm{~d} u \int_{G / T} \mathrm{dV} V_{G / T}(g T) \int_{-\pi}^{\pi} \mathrm{d} \vartheta  \tag{4.4}\\
& \quad\left[\rho_{x}\left(g e^{i \vartheta} g^{-1}\right) \rho_{1}(u)\left(e^{i \vartheta}-e^{-i \vartheta}\right) e^{i k \Psi_{x, y}(u, g T, \vartheta)} s\left(\tilde{\mu}_{g e^{-i \vartheta} g^{-1}}(x), y, k u\right)\right] .
\end{align*}
$$

where

$$
\begin{equation*}
\Psi_{x, y}(u, g T, \vartheta)=u \psi\left(\tilde{\mu}_{g e^{-i \vartheta} g^{-1}}(x), y\right)-\nu \vartheta \tag{4.5}
\end{equation*}
$$

Let $\Im(z)$ denote the imaginary part of $z \in \mathbb{C}$. In view of Corollary 1.3 of [BdMS76], there exists a fixed constant $C_{1}$, depending only on $X$, such that

$$
\begin{equation*}
\Im(\psi(x, y)) \geq C_{1} \operatorname{dist}_{X}(x, y)^{2} . \tag{4.6}
\end{equation*}
$$

Thus, in the range fixed in the hypothesis, we have

$$
\operatorname{dist}_{X}\left(\tilde{\mu}_{g e^{i \vartheta} g_{g^{-1}}}(x), y\right) \geq C k^{\epsilon-1 / 2}
$$

for every $g T \in G / T$ and $e^{i \vartheta} \in T$. In view of (4.5) and (4.6),

$$
\begin{aligned}
& \left|\partial_{u} \Psi_{x, y}(u, g T, \vartheta)\right|=\left|\partial_{u} \psi\left(\tilde{\mu}_{g e^{-i} \vartheta_{g^{-1}}}(x), y\right)\right| \\
& \quad \geq \Im\left(\psi\left(\tilde{\mu}_{g e^{-i \vartheta} \vartheta_{g^{-1}}}(x), y\right)\right) \geq C_{1} \operatorname{dist}_{X}\left(\tilde{\mu}_{g e^{-i} \vartheta_{g^{-1}}}(x), y\right)^{2} \geq C_{1} C^{2} k^{2 \epsilon-1}
\end{aligned}
$$

Eventually, by the identity

$$
-\frac{i}{k} \partial_{u} \psi\left(\tilde{\mu}_{g e^{-i} \vartheta_{g^{-1}}}(x), x\right)^{-1} \frac{\mathrm{~d}}{\mathrm{~d} u} e^{i k \Psi_{x, y}}=e^{i k \Psi_{x, y}}
$$

we can iteratively integrate by parts in $\mathrm{d} u$ in (4.4) and the claim follows.

### 4.2 Proof of Theorem 1.2.2

Proof. Let $\varrho: G \rightarrow[0,+\infty)$ be a smooth bump function supported in a small neighborhood of the origin, and identically equal to 1 on a smaller neighborhood. Let us set $\varrho_{j}(g):=\varrho\left(g g_{j}\right)$. We then have

$$
\begin{align*}
\Pi_{k \nu}(x, x) & =k \nu \cdot \int_{G} \mathrm{~d} V_{G}(g)\left[\overline{\chi_{k \nu}(g)} \Pi\left(\tilde{\mu}_{g^{-1}}(x), x\right)\right]  \tag{4.7}\\
& \sim \sum_{j=0}^{N} k \nu \cdot \int_{G} \mathrm{dV}_{G}(g)\left[\varrho_{j}(g) \cdot \overline{\chi_{k \nu}(g)} \Pi\left(\tilde{\mu}_{g^{-1}}(x), x\right)\right] .
\end{align*}
$$

Let us write $\Pi_{k \nu}(x, x)_{j}$ for the $j$-th summand on the last line of (4.7).
Let $Z_{x}:=G_{x} \cap Z$, where $Z=\left\{ \pm I_{2}\right\}$ is the center of $G$. We shall distinguish two cases, depending on whether $g_{j} \in Z_{x}$ or not. Let us write

$$
\Pi_{k \nu}(x, x) \sim \Pi_{k \nu}(x, x)_{Z_{x}}+\Pi_{k \nu}(x, x)_{G_{x} \backslash Z_{x}}
$$

where

$$
\Pi_{k \nu}(x, x)_{Z_{x}}:=\sum_{g_{j} \in Z_{x}} \Pi_{k \nu}(x, x)_{j}, \quad \Pi_{k \nu}(x, x)_{G_{x} \backslash Z_{x}}:=\sum_{g_{j} \notin Z_{x}} \Pi_{k \nu}(x, x)_{j} .
$$

Let us study $\Pi_{k \nu}(x, x)_{G_{x} \backslash Z_{x}}$. We shall consider the summand in (4.7) with $g_{j} \notin Z_{x}$. Then $g_{j} \neq g_{j}^{-1}$ and thus $G_{x} \backslash Z_{x}$ has even cardinality $b_{x}=2 a_{x}$, and perhaps after renumbering its elements can be arranged in pairs $\left(g_{j}, g_{j}^{-1}\right)$, $j=1, \ldots, a_{x}$.

Let $p: G / T \times T \rightarrow G,(g T, t) \mapsto g t g^{-1}$; since $g_{j}$ is a regular element of $G$, it is a regular value of $p$. For every $j=1, \ldots, a_{x}$, let us set

$$
t_{j}:=h_{m_{x}}^{-1} g_{j} h_{m_{x}}=\left(\begin{array}{cc}
e^{i \vartheta_{j}} & 0 \\
0 & e^{-i \vartheta_{j}}
\end{array}\right) .
$$

We have

$$
p^{-1}\left(g_{j}\right)=\left\{\left(h_{m_{x}} T, t_{j}\right),\left(k_{m_{x}} T, t_{j}^{-1}\right)\right\}, \quad k_{m_{x}}:=h_{m_{x}}\left(\begin{array}{cc}
0 & -1  \tag{4.8}\\
1 & 0
\end{array}\right) .
$$

Definition 4.2.1. If $z \in \mathbb{C}$, let us set

$$
A(z):=i\left(\begin{array}{ll}
0 & z \\
\bar{z} & 0
\end{array}\right) \in \mathfrak{g} .
$$

Then the $\mathbb{R}$-linear map

$$
\eta_{j}: z \in \mathbb{C} \mapsto\left(\operatorname{Ad}_{t_{j}^{-1}}-\operatorname{id}_{\mathfrak{g}}\right)(A(z)) \in \mathfrak{g}
$$

is injective. Therefore, since $\tilde{\mu}$ is locally free at $x$, there is a positive define $2 \times 2$ matrix $C(x ; j)$ such that

$$
\left\|\operatorname{Ad}_{h_{m_{x}}}\left(\eta_{j}(z)\right)_{X}(x)\right\|_{x}^{2}=\frac{1}{2} \cdot Z^{t} C(x ; j) Z \quad(z \in \mathbb{C})
$$

where $Z:=(a, b)^{t} \in \mathbb{R}^{2}$ if $z=a+i b$. Let us define

$$
B(x ; j):=C(x ; j)+4 i \sin \left(2 \vartheta_{j}\right) \cdot \lambda\left(m_{x}\right) I_{2} .
$$

We shall prove the following.
Proposition 4.2.1. Assume that $G_{x} \backslash Z_{x}=\left\{g_{1}, g_{1}^{-1}, \ldots, g_{a_{x}}, g_{a_{x}}^{-1}\right\}$. Then as $k \rightarrow+\infty$ there is an asymptotic expansion

$$
\begin{aligned}
\Pi_{k \nu}(x, x)_{G_{x} \backslash Z_{x}} \sim & 4 \pi \cdot D_{G / T} \cdot\left(\frac{\nu k}{2 \pi \lambda\left(m_{x}\right)}\right)^{d} \\
& \cdot\left[\sum_{j=1}^{a_{x}} \Re\left(\frac{i \sin \left(\vartheta_{j}\right) \cdot e^{-i k \nu \cdot \vartheta_{j}}}{\sqrt{\operatorname{det}(B(x ; j))}}\right)+\sum_{l \geq 1} k^{-l / 2} P_{j l}\left(m_{x}\right)\right] .
\end{aligned}
$$

Proof of Proposition 4.2.1. Let $E: \mathfrak{g} \rightarrow G$ be the exponential map. We shall write the general $t \in T$ in exponential form as

$$
\left(\begin{array}{cc}
e^{i \vartheta} & 0 \\
0 & e^{-i \vartheta}
\end{array}\right)=E(i \vartheta B), \quad B:=\left(\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right)
$$

By the Weyl integration formula and character formulae, we have

$$
\begin{align*}
\Pi_{k \nu}(x, x)_{j}= & \frac{k \nu}{2 \pi} \cdot \int_{G / T} \mathrm{dV} V_{G / T}(g T) \int_{-\pi}^{\pi} \mathrm{d} \vartheta  \tag{4.9}\\
& {\left[\rho_{j}\left(g e^{i \vartheta B} g^{-1}\right) \cdot e^{-i k \nu \cdot \vartheta} \Pi\left(\tilde{\mu}_{g e^{-i \vartheta B} g^{-1}}(x), x\right)\left(e^{i \vartheta}-e^{-i \vartheta}\right)\right] }
\end{align*}
$$

The pulled-back cut-off $\left(g T, e^{i \vartheta}\right) \mapsto \varrho_{j}\left(g e^{i \vartheta B} g^{-1}\right)$ is supported in a small open neighborhood of the pair (4.8). Therefore, we can further split (4.9) as

$$
\Pi_{k \nu}(x, x)_{j}=\Pi_{k \nu}(x, x)_{j 1}+\Pi_{k \nu}(x, x)_{j 2}
$$

where in $\Pi_{k \nu}(x, x)_{j 1}$ (respectively, $\left.\Pi_{k \nu}(x, x)_{j 2}\right)$ integration is over a small neighborhood of $\left(h_{m x} T, t_{j}\right)$ (respectively, $\left(k_{m x} T, t_{j}^{-1}\right)$ ).

Let us consider each $\Pi_{k \nu}(x, x)_{j l}$ separately, starting with $l=1$.
Let us introduce local coordinates on $G / T$ and on $T$. First, for $z \in$ $D(0, \delta) \subset \mathbb{C}$ for some suitably small $\delta>0$, we set

$$
h(z):=h_{m_{x}} E(A(z)),
$$

where $A(z)$ is as in Definition 1.2.2; then the assignment $z \in D(0, \delta) \mapsto$ $h(z) T \in G / T$ is a system of local coordinates on $G / T$ centered at $h_{m_{x}} T$. The Haar measure on $G / T$, expressed in the $z$ coordinates, is $\mathcal{V}_{G / T}(z) \mathrm{d} V_{\mathbb{C}}(z)$, for an appropriate smooth function $\mathcal{V}_{G / T}$. The proof of the following Lemma will be omitted.

Lemma 4.2.1. Let us set $D_{G / T}=\mathcal{V}_{G / T}(0)$. Then $D_{G / T}=2 \pi / V_{3}$, where $V_{3}$ is the surface of $S^{3}$.

Next,

$$
\theta \in(-\delta, \delta) \mapsto t_{j} E(i \theta B) \in T
$$

is a system of local coordinates on T centered at $t_{j}$. Furthermore, since $\left(\tilde{\mu}_{g e^{-i \vartheta B} g^{-1}}(x), x\right)$ is in a small neighborhood of the diagonal in $X \times X$, we may replace $\Pi$ by its representation as an FIO. After performing the rescaling $u \mapsto k u$, we obtain

$$
\begin{align*}
\Pi_{k \nu}(x, x)_{j 1} \sim & \frac{k^{2} \nu}{2 \pi} \cdot e^{-i k \nu \cdot \vartheta_{j}} \cdot \int_{D(0, \delta)} \mathrm{d} \mathrm{~V}_{\mathbb{C}}(z) \int_{-\delta}^{\delta} \mathrm{d} \theta \int_{0}^{+\infty} \mathrm{d} u  \tag{4.10}\\
& {\left[e^{-i k\left[u \psi\left(\tilde{\mu}_{h(z) E(-i \vartheta B) t_{j}^{-1} h(z)^{-1}(x), x}\right)-\nu \cdot \theta\right]}\left(e^{i\left(\vartheta_{j}+\theta\right)}-e^{-i\left(\vartheta_{j}+\theta\right)}\right)\right.} \\
& \left.s_{j 1}\left(\tilde{\mu}_{h(z) E(-i \vartheta B) t_{j}^{-1} h(z)^{-1}}(x), x, k u\right) \mathcal{V}_{G / T}(z)\right] .
\end{align*}
$$

Here, $\mathrm{dV}_{\mathbb{C}}(z)$ is the Lebesgue measure on $\mathbb{C} \cong \mathbb{R}^{2}$, and $s_{j 1}$ denotes the usual amplitude of the representation of $\Pi$ as an FIO, with the above cutoffs incorporated. In addition, by the same argument used before, only a rapidly decreasing contribution is lost if the integrand is multiplied by a bump function in $u$, compactly supported in $(1 / D, D)$, and $\equiv 1$ in $(2 / D, D / 2)$ for some $D \gg 0$ (also implicitly incorporated in the amplitude).

In order to proceed, we need to express the phase more explicitly. We have

$$
\begin{aligned}
h(z) & E(-i \vartheta B) t_{j}^{-1} h(z)^{-1} \\
\quad= & C_{h_{m_{x}}}\left(E(A(z)) E(-i \vartheta B) E\left(-\operatorname{Ad}_{t_{j}^{-1}}(A(z))\right)\right) g_{j}^{-1} \\
\quad= & E\left(-\operatorname{Ad}_{h_{m_{x}}}(\gamma(z, \theta))\right) g_{j}^{-1},
\end{aligned}
$$

where (by the use of Baker-Campbell-Hausdorff)

$$
\gamma(z, \theta)=\gamma_{1}(z, \theta)+\gamma_{2}(z, \theta)+R_{3}(z, \theta)
$$

with

$$
\begin{aligned}
\gamma_{1}(z, \theta) & :=i \theta B+\left(\operatorname{Ad}_{t_{j}^{-1}}-\operatorname{id}_{\mathfrak{g}}\right)(A(z)) \\
\gamma_{2}(z, \theta) & :=-\frac{i}{2}\left[\theta B, A(z)+\operatorname{Ad}_{t_{j}^{-1}}(A(z))\right]+\frac{1}{2}\left[A(z), \operatorname{Ad}_{t_{j}^{-1}}(A(z))\right] .
\end{aligned}
$$

while $R_{j}$ denotes a generic $C^{\infty}$ function vanishing to $j$-th order at the origin. By Corollary 2.2 of [Pao12], we obtain in HLC's

$$
\tilde{\mu}_{h(z) E(-i \vartheta B) t_{j}^{-1} h(z)^{-1}}(x)=\tilde{\mu}_{E\left(-\operatorname{Ad}_{h_{m_{x}}}(\gamma(z, \theta))\right)}(x)=x+(\Theta(z, \theta), V(\theta, z)),
$$

where

$$
\begin{aligned}
\Theta(z, \theta) & :=\left\langle\Phi_{G}\left(m_{x}\right), \operatorname{Ad}_{h_{m_{x}}}(\gamma(z, \theta))\right\rangle+R_{3}(z, \theta) \\
V(\theta, z) & :=\operatorname{Ad}_{h_{m_{x}}}(\gamma(z, \theta))_{M}(m)+R_{2}(z, \theta)
\end{aligned}
$$

By the discussion $\S 3$ of [SZ02] (see especially (65)) we conclude that

$$
\begin{aligned}
u \psi & \left(\tilde{\mu}_{h(z) E(-i \vartheta B) t_{j}^{-1} h(z)^{-1}}(x), x\right)-\nu \theta \\
& =u \psi(x+(\Theta(z, \theta), V(\theta, z)), x)-\nu \theta \\
& =i u \cdot\left[1-e^{i \Theta(z, \theta)}\right]+\frac{i u}{2} \cdot\|V(\theta, z)\|^{2}+u R_{3}(z, \theta) \\
& =u \Theta(z, \theta)+\frac{i u}{2} \cdot\left[\Theta(z, \theta)^{2}+\|V(\theta, z)\|^{2}\right]+u R_{3}(z, \theta) .
\end{aligned}
$$

Let us choose $C>0, \epsilon \in(0,1 / 6)$. Since $p$ is a local diffeomorphism at $\left(h_{m_{x}} T, t_{j}\right)$ and $\tilde{\mu}$ is locally free at $x$, the contribution to the asymptotics
of (4.10) of the locus where $\|(z, \theta)\| \geq C k^{\epsilon-1 / 6}$ is $O\left(k^{-\infty}\right)$. Adopting the rescaling $z \mapsto z / \sqrt{k}, \theta \mapsto \theta / \sqrt{k}$ we can rewrite (4.10) in the following form:

$$
\begin{aligned}
\Pi_{k \nu}(x, x)_{j 1} \sim & \frac{k^{1 / 2} \nu}{2 \pi} \cdot e^{-i k \nu \cdot \vartheta_{j}} \int_{D(0, \delta)} \mathrm{dV}_{\mathbb{C}}(z) \int_{-\delta}^{\delta} \mathrm{d} \theta \int_{0}^{+\infty} \mathrm{d} u \\
& {\left[e^{i k \Psi_{k}(x ; u, \theta, z)} \mathcal{A}_{k}(x ; u, \theta, z)\right] }
\end{aligned}
$$

here we have set

$$
\begin{aligned}
& i k \Psi_{k}(x ; u, \theta, z):= i \sqrt{k}\left[u \cdot\left\langle\Phi_{G}\left(m_{x}\right), \operatorname{Ad}_{h_{m_{x}}}\left(\gamma_{1}(z, \theta)\right)\right\rangle-\theta \nu\right] \\
&+i u \cdot\left\langle\Phi_{G}\left(m_{x}\right), \operatorname{Ad}_{h_{m_{x}}}\left(\gamma_{2}(z, \theta)\right)\right\rangle \\
&-\frac{u}{2}\left\|\operatorname{Ad}_{h_{m_{x}}}\left(\gamma_{1}(z, \theta)\right)_{X}(x)\right\|^{2}+k R_{3}\left(\frac{z}{\sqrt{k}}, \frac{\theta}{\sqrt{k}}\right), \\
& \mathcal{A}_{k}(x ; u, \theta, z):=s_{j 1}\left(\tilde{\mu}_{\left.h(z / \sqrt{k}) E(-i \vartheta / \sqrt{k} B) t_{j}^{-1} h(z / \sqrt{k})^{-1}(x), x, k u\right)}\right. \\
& \varrho\left(k^{-\epsilon}(z, \theta)\right) \cdot\left(e^{i\left(\vartheta_{j}+\theta / \sqrt{k}\right)}-e^{-i\left(\vartheta_{j}+\theta / \sqrt{k}\right)}\right) \cdot \mathcal{V}_{G / T}\left(\frac{z}{\sqrt{k}}\right),
\end{aligned}
$$

with $\varrho$ an appropriate bump function. Integration in $(z, \theta)$ in (4.11) is over a ball of radius $O\left(k^{-\infty}\right)$ centered at the origin.

The second summand in $\gamma_{1}(z, \theta)$ in $\gamma_{2}(z)$ is anti-diagonal. Therefore,

$$
\begin{aligned}
\left\langle\Phi_{G}\left(m_{x}\right), \operatorname{Ad}_{h_{m_{x}}}\left(\gamma_{1}(z, \theta)\right)\right\rangle & =\left\langle\operatorname{Ad}_{h_{m_{x}}^{-1}}\left(\Phi_{G}\left(m_{x}\right)\right), \gamma_{1}(z, \theta)\right\rangle \\
& =\left\langle i \lambda\left(m_{x}\right) B, i \theta B\right\rangle=2 \lambda\left(m_{x}\right) \theta .
\end{aligned}
$$

Let us define

$$
\begin{equation*}
I_{k}(x, z):=\int_{-\infty}^{+\infty} \mathrm{d} \theta \int_{0}^{+\infty} \mathrm{d} u\left[e^{i \sqrt{k} \Psi_{k}(u, \theta)} \mathcal{A}_{k}(x ; u, \theta, z)\right] \tag{4.12}
\end{equation*}
$$

where now

$$
\begin{aligned}
\Psi_{k}(u, \theta) & :=\theta\left[2 \lambda\left(m_{x}\right) \cdot u-\nu\right] \\
\mathcal{A}_{k}(x ; u, \theta, z) & :=e^{\mathcal{E}_{k}(u, \theta, z)} \cdot A_{k}(x ; u, \theta, z),
\end{aligned}
$$

with

$$
\begin{align*}
\mathcal{E}(u, \theta, z):= & i u \cdot\left\langle\Phi_{G}\left(m_{x}\right), \operatorname{Ad}_{h_{m_{x}}}\left(\gamma_{2}(z, \theta)\right)\right\rangle  \tag{4.13}\\
& -\frac{u}{2} \cdot\left\|\operatorname{Ad}_{h_{m_{x}}}\left(\gamma_{1}(z, \theta)_{X}(x)\right)\right\|_{x}^{2}+k R_{3}\left(\frac{z}{\sqrt{k}}, \frac{\theta}{\sqrt{k}}\right) .
\end{align*}
$$

Since $\tilde{\mu}$ is locally free at $x$, on the support of the integrand

$$
\Re\left(\mathcal{E}_{k}(z, \theta)\right) \leq-D^{\prime} \cdot\left(|z|^{2}+|\theta|^{2}\right)
$$

for some positive constant $D>0$.
Then we can rewrite (4.11) in the following form:

$$
\begin{equation*}
\Pi_{k \nu}(x, x)_{j 1} \sim \frac{k^{1 / 2} \nu}{2 \pi} \cdot e^{-i k \nu \cdot \vartheta_{j}} \cdot \int_{\mathbb{C}} \mathrm{d} V_{\mathbb{C}}(z)\left[I_{k}(x, z)\right] \tag{4.14}
\end{equation*}
$$

The following is straightforward.
Lemma 4.2.2. $\Psi_{k}$ has a unique critical point, which is non-degenerate and given by $\left(u_{0}, \theta_{0}\right)=\left(\nu /\left(2 \lambda\left(m_{x}\right)\right), 0\right)$; we have $\Psi_{k}\left(u_{0}, \theta_{0}\right)=0$. The Hessian matrix has determinant $-4 \lambda\left(m_{x}\right)^{2}$ and vanishing signature.

We can apply the Stationary Phase Lemma to determine the asymptotic expansion of (4.12). By a few computations we get

$$
\begin{aligned}
& \gamma_{1}(z, 0)=i\left(\begin{array}{cc}
0 & \left(e^{-2 i \vartheta_{j}}-1\right) \cdot z \\
\left(e^{2 i \vartheta_{j}}-1\right) \cdot \bar{z} & 0
\end{array}\right), \\
& \gamma_{2}(z, 0)=-i \cdot|z|^{2} \cdot \sin \left(2 \vartheta_{j}\right) B .
\end{aligned}
$$

If $z=a+i b$ with $a, b \in \mathbb{R}$, let $Z=(a, b)^{t} \in \mathbb{R}^{2}$ be the corresponding vector; thus $|z|=\|Z\|$. Then

$$
\left\|\operatorname{Ad}_{h_{m_{x}}}\left(\gamma_{1}(z, 0)\right)_{X}(x)\right\|_{x}^{2}=\frac{1}{2} \cdot Z^{t} C^{t}(x, j) Z
$$

where $C(x, j)$ is as in Definition 1.2.2. From (4.13) we conclude

$$
\mathcal{E}_{k}\left(u_{0}, 0, z\right)=-\frac{u_{0}}{2} Z^{t} B(x, j) Z+k R_{3}\left(\frac{z}{\sqrt{k}}\right)
$$

where $B(x, j)$ is also as in Definition 1.2.2.

On the other hand, we have (noting that $\sin \left(\vartheta_{j}\right) \neq 0$ as $g_{j} \notin Z_{x}$ )

$$
e^{i\left(\vartheta_{j}+\theta / \sqrt{k}\right)}-e^{-i\left(\vartheta_{j}+\theta / \sqrt{k}\right)}=2 i \sin \left(\vartheta_{j}\right) \cdot\left[1+\sum_{j \geq 1} k^{-l / 2} a_{j l}(\theta)\right] .
$$

Applying the Stationary Phase Lemma, we obtain that as $k \rightarrow+\infty$

$$
\begin{align*}
I_{k}(x, z) \sim & D_{G / T}\left(\frac{k u_{0}}{\pi}\right)^{d} \cdot \frac{\pi}{\sqrt{k}} \cdot \frac{2 i \sin \left(\vartheta_{j}\right)}{\lambda\left(m_{x}\right)} \cdot e^{-\frac{u_{0}}{2} Z^{t} A(x, j) Z}  \tag{4.15}\\
& \cdot\left[1+\sum_{l \geq 1} k^{-l / 2} R_{j l}\left(m_{x}, Z\right)\right]
\end{align*}
$$

where $R_{j l}\left(m_{x}, Z\right)$ is polynomial in $Z$ of degree $\leq 3 l$. Inserting (4.15) in (4.14), we obtain

$$
\begin{array}{r}
\Pi_{k \nu}(x, x)_{j 1} \sim 4 \pi \cdot D_{G / T} \cdot \frac{\sin \left(\vartheta_{j}\right) \cdot e^{-i k \nu \cdot \vartheta_{j}}}{\sqrt{\operatorname{det} B(x, j)}} \cdot\left(\frac{\nu k}{2 \pi \cdot \lambda\left(m_{x}\right)}\right)^{d}  \tag{4.16}\\
\cdot\left[1+\sum_{l \geq 1} k^{-l / 2} P_{j l}\left(m_{x}\right)\right] .
\end{array}
$$

Let us remark that $g_{j} \neq g_{j}^{-1}$ since $g_{j} \neq \pm I_{2}$; summing the contributions (4.16) corresponding to $g_{j}$ and $g_{j^{\prime}}=g_{j}^{-1}$, we obtain

$$
\begin{aligned}
\Pi_{k \nu}(x, x)_{j 1}+\Pi_{k \nu}(x, x)_{j^{\prime} 1}= & 8 \pi \cdot \Re\left(\frac{i \sin \left(\vartheta_{j}\right) \cdot e^{-i k \nu \vartheta_{j}}}{\sqrt{\operatorname{det}(A(x, j))}}\right) \cdot\left(\frac{\nu k}{2 \pi \cdot \lambda\left(m_{x}\right)}\right)^{d} \\
& \cdot\left[D_{G / T}+\sum_{l \geq 1} k^{-l / 2} P_{j l}\left(m_{x}\right)\right]
\end{aligned}
$$

To deal with $\Pi_{k \nu}(x, x)_{j 2}$, in view of (4.8) we need only go over the previous computations replacing $h_{m_{x}}$ with $k_{m_{x}}$, and $t_{j}$ with $t_{j}^{-1}$. In the analogue of (4.12), in place of the phase $\Psi_{x}$, we obtain

$$
\Psi_{x}^{\prime}(u, \theta)=-\theta \cdot\left[2 \lambda\left(m_{x}\right) \cdot u+\nu\right]
$$

so that $\partial_{\theta} \Psi_{x}^{\prime}(u, \theta)=-\left[2 \lambda\left(m_{x}\right) \cdot u+\nu\right] \leq-\nu$ on the domain of integration. Thus $\Pi_{k \nu}(x, x)_{j 2}=O\left(k^{-\infty}\right)$.

The proof of Proposition 4.2.1 is complete.

### 4.3 Proof of Theorem 1.2.3

Proof of Theorem 1.2.3. We have

$$
\Pi_{k \nu}\left(x_{1 k}, x_{2 k}\right)=k \nu \int_{G} \mathrm{dV}_{G}(g)\left[\overline{\chi_{k \nu}(g)} \Pi\left(\tilde{\mu}_{g^{-1}}\left(x_{1 k}\right), x_{2 k}\right)\right] .
$$

Only a rapidly decreasing contribution to the asymptotic is lost, if integration is restricted to a small neighborhood of $G_{x}$. Thus we may multiply the integrand by a cut-off function $\varrho \in \mathcal{C}^{\infty}(G)$ supported in a small neighborhood of $G_{x}$ and $\equiv 1$ in a slightly smaller neighborhood, without altering the asymptotics. We may assume that $\varrho$ is invariant under conjugation; thus we shall write $\varrho=\varrho(t)$, or $\varrho(\vartheta)$ working in coordinates.

By the Weyl integration formula and character formula,

$$
\begin{aligned}
\Pi_{k \nu}\left(x_{1 k}, x_{2 k}\right) \sim & k \nu \int_{T} \mathrm{~d} V_{T}(t) \int_{G / T} \mathrm{~d} \mathrm{~V}_{G / T}(g T) \\
= & \quad\left[t^{-k \nu}\left(t-t^{-1}\right) \varrho(t) \Pi\left(\tilde{\mu}_{g t-1} g^{-1}\left(x_{1 k}, x_{2 k}\right)\right)\right] \\
2 \pi & \int_{-\pi / 2}^{3 \pi / 2} \mathrm{~d} \vartheta \int_{G / T} \mathrm{~d} V_{G / T}(g T) \\
& \quad\left[e^{-i k \nu \vartheta}\left(e^{i \vartheta}-e^{-i \vartheta}\right) \varrho(\vartheta) \Pi\left(\tilde{\mu}_{g e^{-i \vartheta} B_{g} g^{-1}}\left(x_{1 k}\right), x_{2 k}\right)\right] .
\end{aligned}
$$

On the support of $\varrho,\left(\tilde{\mu}_{g e^{-i \vartheta B} g^{-1}}\left(x_{1 k}\right), x_{2 k}\right)$ lies in a small neighborhood of the diagonal; hence we may replace $\Pi$ by its representation as an FIO, without changing the asymptotics. Therefore,

$$
\begin{aligned}
\Pi_{k \nu}\left(x_{1 k}, x_{2 k}\right) \sim & \frac{k^{2} \nu}{2 \pi} \int_{0}^{+\infty} \mathrm{d} u \int_{-\pi / 2}^{3 \pi / 2} \mathrm{~d} \vartheta \int_{G / T} \mathrm{~d} V_{G / T}(g T)\left[\rho(\vartheta)\left(e^{i \vartheta}-e^{-i \vartheta}\right)\right. \\
& \left.e^{i k\left[u \psi\left(\tilde{\mu}_{g e^{-i \vartheta}} B_{g^{-1}}\left(x_{1 k}\right), x_{2 k}\right)-\nu \vartheta\right]} s\left(\tilde{\mu}_{g e^{-i \vartheta B} g^{-1}}\left(x_{1 k}\right), x_{2 k}, k u\right)\right],
\end{aligned}
$$

where we have performing the rescaling $u \mapsto k u$.
The same arguments of Lemmas 4.1.1 and 4.1.4 apply here with minor modifications. In particular, we have

Lemma 4.3.1. Let $D \gg 0$ and let $\varrho_{1} \in \mathcal{C}_{c}(\mathbb{R})$ be $\geq 0$, supported in $(1 / D, D)$, and $\equiv 1$ on $(2 / D, D / 2)$. Then only a rapidly decreasing contribution to the asymptotics is lost, if the integrand is multiplied by $\varrho_{1}(u)$.

We obtain

$$
\begin{array}{r}
\Pi_{k \nu}\left(x_{1 k}, x_{2 k}\right) \sim \frac{k^{2} \nu}{2 \pi} \int_{0}^{+\infty} \mathrm{d} u \int_{-\pi / 2}^{3 \pi / 2} \mathrm{~d} \vartheta \int_{G / T} \mathrm{dV}  \tag{4.17}\\
G / T \\
(g T)\left[\varrho_{1}(u)\left(e^{i \vartheta}-e^{-i \vartheta}\right)\right. \\
\left.\rho(\vartheta) e^{i k\left[u \psi\left(\tilde{\mu}_{g e^{-i \vartheta} B_{g-1}}\left(x_{1 k}\right), x_{2 k}\right)-\nu \vartheta\right]} s\left(\tilde{\mu}_{g e^{-i \vartheta B} g^{-1}}\left(x_{1 k}\right), x_{2 k}, k u\right)\right]
\end{array}
$$

We can write $\varrho(\vartheta)=\varrho_{0}(\vartheta)+\varrho_{\pi}(\vartheta)$, where $\varrho_{0}(\vartheta)$ is supported in a small neighborhood of 0 , and identically equal to 1 in a smaller neighborhood; on the other hand, $\varrho_{\pi}(\vartheta)$ is supported in a small neighborhood of $\pi$, and in fact it vanishes identically if $-I_{2} \notin G_{x}$, while it is identically equal to 1 on a smaller neighborhood if $I_{2} \in G_{x}$.

Inserting the latter identity in (4.17), we shall accordingly write

$$
\Pi_{k \nu}\left(x_{1 k}, x_{2 k}\right) \sim \Pi_{k \nu}\left(x_{1 k}, x_{2 k}\right)_{0}+\Pi_{k \nu}\left(x_{1 k}, x_{2 k}\right)_{\pi},
$$

and examine the two summands separately.
First, let us consider the asymptotics of $\Pi_{k \nu}\left(x_{1 k}, x_{2 k}\right)_{0}$. We shall prove the following.

Proposition 4.3.1. Under the assumptions of Theorem 1.2.3, as $k \rightarrow+\infty$ we have

$$
\begin{aligned}
\Pi_{k \nu}\left(x_{1 k}, x_{2 k}\right)_{0} \sim & \frac{1}{2 \lambda\left(m_{x}\right)} \cdot\left(\frac{\nu k}{2 \pi \lambda\left(m_{x}\right)}\right)^{d} \cdot e^{u_{0} \psi_{2}\left(\mathbf{v}_{1}, \mathbf{v}_{2}\right)} \\
& \cdot\left[1+\sum_{j \geq 1}^{+\infty} k^{-j / 2} A_{j}\left(x, \mathbf{v}_{1}, \mathbf{v}_{2}\right)\right]
\end{aligned}
$$

where $A_{j}(x, \cdot, \cdot)$ is a polynomial of degree $\leq 3 j$ and parity $(-1)^{j}$.
Proof of Proposition 4.3.1. Here $\varrho(\vartheta)$ has been replaced by $\varrho_{0}(\vartheta)$, therefore integration in $\mathrm{d} \vartheta$ is restricted to $(-2 \delta, 2 \delta)$. Let us fix $C_{1}>0, \epsilon_{1} \in(0,1 / 6)$. Iteratively integrating by parts in $\mathrm{d} u$, we conclude that the locus where $|\vartheta|>C_{1} k^{\epsilon_{1}-1 / 2}$ contributes negligibly to the asymptotics of $\Pi_{k \nu}\left(x_{1 k}, x_{2 k}\right)_{0}$. More precisely, we have the following.
Lemma 4.3.2. Suppose that $\varrho_{2} \in \mathcal{C}_{c}(\mathbb{R})$ is $\geq 0$, supported in $(-2,2)$, and $\equiv 1$ on $(-1,1)$. Then the asymptotics of $\Pi_{k \nu}\left(x_{1 k}, x_{2 k}\right)_{0}$ are unchanged, if the integrand is multiplied by $\varrho_{2}\left(k^{1 / 2-\epsilon_{1}} \vartheta\right)$.

Applying the rescaling $\vartheta \mapsto \vartheta / \sqrt{k}$, we recover

$$
\begin{aligned}
& \Pi_{k \nu}\left(x_{1 k}, x_{2 k}\right)_{0} \\
& \sim \frac{k^{3 / 2} \nu}{2 \pi} \int_{0}^{+\infty} \mathrm{d} u \int_{-\infty}^{+\infty} \mathrm{d} \vartheta \int_{G / T} \mathrm{dV}_{G / T}(g T)\left[\rho_{2}\left(k^{-\epsilon_{1}} \vartheta\right) \varrho_{1}(u) e^{i k \Psi_{k}}\right. \\
& \left.\quad\left(e^{i \vartheta / \sqrt{k}}-e^{-i \vartheta / \sqrt{k}}\right) s\left(\tilde{\mu}_{g e^{-i \vartheta B / \sqrt{k}} g^{-1}}\left(x_{1 k}\right), x_{2 k}, k u\right)\right],
\end{aligned}
$$

where

$$
\Psi_{k}\left(u, v_{1}, v_{2}, \vartheta, g T\right):=\psi\left(\tilde{\mu}_{g e^{-i \vartheta B / \sqrt{k}} g^{-1}}\left(x_{1 k}\right), x_{2 k}\right)-\frac{\vartheta}{\sqrt{k}} \nu
$$

Integration in $\mathrm{d} \vartheta$ is effectively over an interval of length $4 k^{\epsilon_{1}}$ centered at the origin.

The next step is to make $\Psi_{k}$ more explicit. By Corollary 2.2 of [Pao12], with $m_{x}=\pi(x)$ we have

$$
\begin{aligned}
& \tilde{\mu}_{g e^{-i \vartheta B / \sqrt{k}} g^{-1}}\left(x_{1 k}\right)=\tilde{\mu}_{e^{-\vartheta \operatorname{Ad} g(\beta) / \sqrt{k}}}\left(x_{1 k}\right) \\
& \quad x+\left(\Theta_{k}\left(v_{1}, \vartheta, g T\right), \frac{1}{\sqrt{k}} V\left(v_{1}, \vartheta, g T\right)+R_{2}\left(\frac{1}{\sqrt{k}} \vartheta, \frac{1}{\sqrt{k}} \mathbf{v}_{1}\right)\right),
\end{aligned}
$$

where (for appropriate $R_{3}$ and $R_{2}$ )

$$
\begin{aligned}
\Theta_{k}\left(v_{1}, \vartheta, g T\right):= & \frac{1}{\sqrt{k}}\left[\theta_{1}+\vartheta \cdot\left\langle\Phi_{G}\left(m_{x}\right), \operatorname{Ad}_{g}(\beta)\right\rangle\right] \\
& +\frac{1}{k} \vartheta \cdot \omega_{m_{x}}\left(\operatorname{Ad}_{g}(\beta)_{M}(m), \mathbf{v}_{1}\right)+R_{3}\left(\frac{1}{\sqrt{k}} \vartheta, \frac{1}{\sqrt{k}} \mathbf{v}_{1}\right)
\end{aligned}
$$

and

$$
V\left(\mathbf{v}_{1}, \vartheta, g T\right):=\mathbf{v}_{1}-\vartheta \operatorname{Ad}_{g}(\beta)_{M}\left(m_{x}\right)
$$

Let us set

$$
\tilde{\Theta}_{k}\left(v_{1}, v_{2}, \vartheta, g T\right):=\frac{1}{\sqrt{k}} A+\frac{1}{k} B+R_{3}\left(\frac{1}{\sqrt{k}} \vartheta, \frac{1}{\sqrt{k}} \mathbf{v}_{1}\right)
$$

where

$$
A=A\left(v_{1}, v_{2}, \vartheta, g T\right):=\vartheta \cdot\left\langle\Phi_{G}\left(m_{x}\right), \operatorname{Ad}_{g}(\beta)\right\rangle
$$

and

$$
B=B\left(v_{1}, \vartheta, g T\right):=\vartheta \cdot \omega_{m_{x}}\left(\operatorname{Ad}_{g}(\beta)_{M}(m), \mathbf{v}_{1}\right)
$$

Then, in view of the discussion of $\S 3$ of [SZ02] (see especially (65)), we conclude that

$$
\begin{aligned}
\Psi_{k}\left(u, v_{1}, v_{2}, \vartheta, g T\right)= & i u\left[1-e^{\tilde{\Theta}_{k}}\right]-\frac{\vartheta}{\sqrt{k}} \nu-i \frac{u}{k} \psi_{2}\left(V, \mathbf{v}_{2}\right) \\
& +R_{3}\left(\frac{1}{\sqrt{k}}\left(\vartheta, v_{1}, v_{2}\right)\right)
\end{aligned}
$$

By a few computations, this leads to the following.
Lemma 4.3.3. We have

$$
\begin{aligned}
\Psi_{k}\left(u, v_{1}, v_{2}, \vartheta, g T\right)= & \frac{1}{\sqrt{k}} \mathcal{G}_{\theta_{1}, \theta_{2}}(u, \vartheta, g T)+\frac{1}{k} \mathcal{D}\left(u, v_{1}, v_{2}, \vartheta, g T\right) \\
& +R_{3}\left(\frac{1}{\sqrt{k}}\left(\vartheta, v_{1}, v_{2}\right)\right) .
\end{aligned}
$$

where

$$
\mathcal{G}_{\theta_{1}, \theta_{2}}(u, \vartheta, g T)=u A-\vartheta \nu
$$

and

$$
\mathcal{D}\left(u, v_{1}, v_{2}, \vartheta, g T\right)=u\left[B+i\left(\frac{1}{2} A^{2}-\psi_{2}\left(V, \mathbf{v}_{2}\right)\right)\right]
$$

Thus, we conclude that

$$
\begin{align*}
\Pi_{k \nu}\left(x_{1 k}, x_{2 k}\right)_{0} \sim & \frac{k^{3 / 2} \nu}{2 \pi} \int_{0}^{+\infty} \mathrm{d} u \int_{-\infty}^{+\infty} \mathrm{d} \vartheta \int_{G / T} \mathrm{~d} V_{G / T}(g T)  \tag{4.18}\\
& {\left[e^{i \sqrt{k} \mathcal{G}_{\theta_{1}, \theta_{2}}(u, \vartheta, g T)} e^{i u B+u\left[\psi_{2}\left(V, \mathbf{v}_{2}\right)-\frac{1}{2} A^{2}\right]} \rho_{2}\left(k^{-\epsilon_{1}} \vartheta\right)\right.} \\
& \left.\varrho_{1}(u)\left(e^{i \vartheta / \sqrt{k}}-e^{-i \vartheta / \sqrt{k}}\right) s\left(\tilde{\mu}_{g e^{-i \vartheta B / \sqrt{k}} g^{-1}}\left(x_{1 k}\right), x_{2 k}, k u\right)\right] .
\end{align*}
$$

In order to make (4.18) yet more explicit, let us make recourse to the coordinates $(\theta, \delta)$ on $G / T$ discussed in section $\S 3.2 .1$. Thus we shall make the replacement

$$
\int_{G / T} \mathrm{~d} V_{G / T}(g T) \mapsto \frac{1}{2 \pi} \int_{0}^{\pi / 2} \mathrm{~d} \theta \int_{-\pi}^{\pi} \mathrm{d} \delta \sin (2 \theta)
$$

Furthermore, let $h_{m} T \in G / T$ be as in the introduction, we shall operate the change of variable $g T \mapsto h_{m} g T$ in $G / T$, and write $g$ as in section $\S 3.2 .1$ with $\alpha=\cos (\theta) e^{i \delta}$ and $\beta=\sin (\theta)$. Then

$$
\begin{align*}
\mathcal{G}_{\theta_{1}, \theta_{2}}\left(u, \vartheta, h_{m} g T\right) & =u \cdot\left[\vartheta \cdot\left\langle i g^{-1}\left(\begin{array}{cc}
\lambda\left(m_{x}\right) & 0 \\
0 & -\lambda\left(m_{x}\right)
\end{array}\right) g, \beta\right\rangle\right]-\vartheta \nu  \tag{4.19}\\
& =u \cdot\left[2 \vartheta \cdot \cos (2 \theta) \cdot \lambda\left(m_{x}\right)\right]-\vartheta \nu
\end{align*}
$$

Let $\mathcal{G}_{\theta_{1}, \theta_{2}}^{\prime}(u, \vartheta, \theta)$ denote the expression on the last line of (4.19). We can rewrite (4.18) in the following form

$$
\begin{align*}
& \Pi_{k \nu}\left(x_{1 k}, x_{2 k}\right)_{0} \sim \frac{k^{3 / 2} \nu}{4 \pi^{2}} \int_{0}^{+\infty} \mathrm{d} u \int_{-\infty}^{+\infty} \mathrm{d} \vartheta \int_{0}^{\pi / 2} \mathrm{~d} \theta \int_{-\pi}^{\pi} \mathrm{d} \delta  \tag{4.20}\\
& {\left[e^{i \sqrt{k} \mathcal{G}_{\theta_{1}, \theta_{2}}(u, \vartheta, g T)} e^{i u B+u\left[\psi_{2}\left(V, \mathbf{v}_{2}\right)-\frac{1}{2} A^{2}\right]} \rho_{2}\left(k^{-\epsilon_{1}} \vartheta\right)\right.} \\
& \left.\quad \varrho_{1}(u)\left(e^{i \vartheta / \sqrt{k}}-e^{-i \vartheta / \sqrt{k}}\right) s\left(\tilde{\mu}_{g e^{-i \vartheta B / \sqrt{k}} g^{-1}}\left(x_{1 k}\right), x_{2 k}, k u\right) \sin (2 \theta)\right],
\end{align*}
$$

where (with abuse of notation) $g=g(\theta, \delta)$ and $A=A(\theta, \delta), B=B(\theta, \delta)$ with the obvious change of variables.

By the change of integration variable $t=\cos (2 \theta)$, we can further reformulate (4.20) as follows. With some abuse of notation, let us write $g T=g(t, \delta) T$ and set

$$
\begin{aligned}
\Gamma_{\theta_{1}, \theta_{2}}(u, \vartheta, t) & :=\mathcal{G}_{\theta_{1}, \theta_{2}}^{\prime}(u, \vartheta, \theta) \\
& =u \cdot\left[2 \vartheta \cdot t \cdot \lambda\left(m_{x}\right)\right]-\vartheta \nu .
\end{aligned}
$$

Then

$$
\begin{align*}
& \Pi_{k \nu}\left(x_{1 k}, x_{2 k}\right)_{0} \sim \frac{1}{2} \cdot \frac{k^{3 / 2} \nu}{(2 \pi)^{2}} \int_{0}^{+\infty} \mathrm{d} u \int_{-\infty}^{+\infty} \mathrm{d} \vartheta \int_{-1}^{1} \mathrm{~d} t \int_{-\pi}^{\pi} \mathrm{d} \delta  \tag{4.21}\\
& {\left[e^{i \sqrt{k} \Gamma_{\theta_{1}, \theta_{2}}(u, \vartheta, t)} e^{i u B_{t}+u\left[\psi_{2}\left(V, \mathbf{v}_{2}\right)-\frac{1}{2} A_{t}^{2}\right]} \rho_{2}\left(k^{-\epsilon_{1}} \vartheta\right)\right.} \\
& \left.\quad \varrho_{1}(u)\left(e^{i \vartheta / \sqrt{k}}-e^{-i \vartheta / \sqrt{k}}\right) s\left(\tilde{\mu}_{g e^{-i \vartheta B / \sqrt{k}} g^{-1}}\left(x_{1 k}\right), x_{2 k}, k u\right) \sin (2 \theta)\right]
\end{align*}
$$

we have denoted by $A_{t}, B_{t}$ the functions

$$
A_{t}\left(v_{1}, v_{2}, \vartheta, \delta\right)=A_{t}\left(v_{1}, v_{2}, \vartheta, g(t, \delta) T\right), B_{t}\left(v_{1}, \vartheta, \delta\right)=B_{t}\left(v_{1}, \vartheta, g(t, \delta) T\right)
$$

and similarly for $V_{t}$.
Let us remark that

$$
\begin{aligned}
e^{i \vartheta / \sqrt{k}}-e^{-i \vartheta / \sqrt{k}} & =\frac{2 i}{\sqrt{k}} \vartheta \sum_{j=0}^{+\infty} \frac{(-1)^{j}}{(2 j+1)!} \frac{\vartheta^{2 j}}{k^{j}} \\
& =\frac{2 i}{\sqrt{k}} \vartheta\left[1+R_{2}\left(\frac{\vartheta}{\sqrt{k}}\right)\right] .
\end{aligned}
$$

Furthermore, working in HLC's, Taylor expansion yields an asymptotic expansion

$$
s\left(\tilde{\mu}_{g e^{-i \vartheta B / \sqrt{k}} g^{-1}}\left(x_{1 k}\right), x_{2 k}, k u\right) \sim\left(\frac{k u}{\pi}\right)^{d}\left[1+R_{1}\left(\frac{\vartheta}{\sqrt{k}}\right)\right] .
$$

It follows that (4.21) is given by an asymptotic expansion in descending halfinteger powers of $k$; with a few computations, one sees that the dominant term is to be extracted from

$$
\begin{align*}
& I_{v_{1}, v_{2}}(k):=\frac{\nu}{8 \pi} \cdot\left(\frac{k}{\pi}\right)^{d+1} \int_{0}^{+\infty} \mathrm{d} u \int_{-\infty}^{+\infty} \mathrm{d} \vartheta \int_{-1}^{1} \mathrm{~d} t \int_{-\pi}^{\pi} \mathrm{d} \delta  \tag{4.22}\\
& {\left[e^{i \sqrt{k} \Gamma_{\theta_{1}, \theta_{2}}(u, \vartheta, t)} \cdot(2 i \vartheta u)\right.} \\
&\left.\quad e^{i u B_{t}+u\left[\psi_{2}\left(V, \mathbf{v}_{2}\right)-\frac{1}{2} A_{t}^{2}\right]} \varrho_{1}(u) \rho_{2}\left(k^{-\epsilon_{1}} \vartheta\right) u^{d-1}\right] ;
\end{align*}
$$

the latter may in turn be rewritten

$$
\begin{aligned}
I_{v_{1}, v_{2}}(k)= & \frac{\nu}{8 \pi} \cdot\left(\frac{k}{\pi}\right)^{d+1} \frac{1}{\sqrt{k} \cdot \lambda\left(m_{x}\right)} \\
& \cdot \int_{0}^{+\infty} \mathrm{d} u \int_{-\infty}^{+\infty} \mathrm{d} \vartheta \int_{-1}^{1} \mathrm{~d} t \int_{-\pi}^{\pi} \mathrm{d} \delta\left[\partial_{t}\left(e^{i \sqrt{k} \Gamma_{\theta_{1}, \theta_{2}}(u, \vartheta, t)}\right)\right. \\
& \left.e^{i u B_{t}+u\left[\psi_{2}\left(V, \mathbf{v}_{2}\right)-\frac{1}{2} A_{t}^{2}\right]} \varrho_{1}(u) \rho_{2}\left(k^{-\epsilon_{1}} \vartheta\right) u^{d-1}\right] .
\end{aligned}
$$

Integrating by parts in $\mathrm{d} t$, we obtain

$$
I_{v_{1}, v_{2}}(k)=\frac{\nu}{8 \pi} \cdot\left(\frac{k}{\pi}\right)^{d+1} \frac{1}{\sqrt{k} \cdot \lambda\left(m_{x}\right)} \cdot\left[J_{v_{1}, v_{2}}^{\prime}(k)-J_{v_{1}, v_{2}}^{\prime \prime}(k)-J_{v_{1}, v_{2}}^{\prime \prime \prime}(k)\right]
$$

where

$$
\begin{aligned}
J_{v_{1}, v_{2}}^{\prime}(k):= & \int_{0}^{+\infty} \mathrm{d} u \int_{-\infty}^{+\infty} \mathrm{d} \vartheta \int_{-\pi}^{\pi} \mathrm{d} \delta\left[e^{i \sqrt{k} \Gamma_{\theta_{1}, \theta_{2}}(u, \vartheta, 1)}\right. \\
& \left.e^{i u B_{1}+u\left[\psi_{2}\left(V, \mathbf{v}_{2}\right)-\frac{1}{2} A_{1}^{2}\right]} \varrho_{1}(u) \rho_{2}\left(k^{-\epsilon_{1}} \vartheta\right) u^{d-1}\right] \\
J_{v_{1}, v_{2}}^{\prime \prime}(k):= & \int_{0}^{+\infty} \mathrm{d} u \int_{-\infty}^{+\infty} \mathrm{d} \vartheta \int_{-\pi}^{\pi} \mathrm{d} \delta\left[e^{i \sqrt{k} \Gamma_{\theta_{1}, \theta_{2}}(u, \vartheta,-1)}\right. \\
& \left.e^{i u B_{-1}+u\left[\psi_{2}\left(V, \mathbf{v}_{2}\right)-\frac{1}{2} A_{-1}^{2}\right]} \varrho_{1}(u) \rho_{2}\left(k^{-\epsilon_{1}} \vartheta\right) u^{d-1}\right],
\end{aligned}
$$

and

$$
\begin{aligned}
J_{v_{1}, v_{2}}^{\prime \prime \prime}(k):= & \int_{-1}^{1} \mathrm{~d} t \int_{0}^{+\infty} \mathrm{d} u \int_{-\infty}^{+\infty} \mathrm{d} \vartheta \int_{-\pi}^{\pi} \mathrm{d} \delta\left[e^{i \sqrt{k} \Gamma_{\theta_{1}, \theta_{2}}(u, \vartheta, t)}\right. \\
& \left.\cdot \partial_{t}\left(e^{i u B_{t}+u\left[\psi_{2}\left(V, \mathbf{v}_{2}\right)-\frac{1}{2} A_{t}^{2}\right]}\right) \varrho_{1}(u) \rho_{2}\left(k^{-\epsilon_{1}} \vartheta\right) u^{d-1}\right] .
\end{aligned}
$$

Let us estimate the three summands separately.
Lemma 4.3.4. As $k \rightarrow+\infty$, there is an asymptotic expansion of the form

$$
\begin{aligned}
& J_{v_{1}, v_{2}}^{\prime}(k) \sim \frac{2 \pi^{2}}{\lambda\left(m_{x}\right) \sqrt{k}} \cdot e^{u_{0} \psi_{2}\left(\mathbf{v}_{1}, \mathbf{v}_{2}\right)} \cdot\left(\frac{\nu}{2 \lambda\left(m_{x}\right)}\right)^{d-1} \\
& \cdot\left[1+\sum_{l=1}^{+\infty} k^{-l / 2} a_{l}\left(m_{x} ; v_{1}, v_{2}\right)\right]
\end{aligned}
$$

where $a_{l}\left(m_{x} ; \cdot, \cdot\right)$ is a polynomial of degree $\leq 3 j$, whose coefficient are smooth function on $M$.

Proof of Lemma 4.3.4. Let us view $J_{v_{1}, v_{2}}^{\prime}(k)$ as an oscillatory integral in the parameter $\sqrt{k}$, with real phase $\Gamma_{\theta_{1}, \theta_{2}}(1 ; u, \vartheta)$

$$
\begin{aligned}
\Gamma_{\theta_{1}, \theta_{2}}(1 ; u, \vartheta) & =u \cdot\left[2 \lambda\left(m_{x}\right) \cdot \vartheta\right]-\vartheta \nu \\
& =u \cdot\left[\left\langle\Phi_{G}\left(m_{x}\right), \operatorname{Ad}_{h_{m_{x}}}(\beta)\right\rangle \cdot \vartheta\right]-\vartheta \nu
\end{aligned}
$$

and amplitude

$$
e^{i u B_{1}+u\left[\psi_{2}\left(V, \mathbf{v}_{2}\right)-\frac{1}{2} A_{1}^{2}\right]} \varrho_{1}(u) \rho_{2}\left(k^{-\epsilon_{1}} \vartheta\right) u^{d-1} .
$$

Explicitly, the exponent is

$$
\begin{aligned}
& \mathcal{E}\left(u, \vartheta, v_{1}, v_{2}\right):=i u B_{1}+u\left[\psi_{2}\left(V, \mathbf{v}_{2}\right)-\frac{1}{2} A_{1}^{2}\right] \\
& =u\left[-i \omega_{m}\left(\mathbf{v}_{1}, \mathbf{v}_{2}\right)+i \vartheta \omega_{m_{x}}\left(\operatorname{Ad}_{h_{m_{x}}}(\beta)_{M}\left(m_{x}\right), \mathbf{v}_{1}+\mathbf{v}_{2}\right)\right. \\
& \left.\quad-\frac{1}{2}\left\|\left(v_{1}-v_{2}\right)-\vartheta \operatorname{Ad}_{h_{m_{x}}}(\beta)_{X}(x)\right\|^{2}\right] .
\end{aligned}
$$

Under the hypothesis of the Theorem, therefore, $\Re\left(\mathcal{E}_{1}\right) \leq-C^{\prime} \vartheta^{2}+C^{\prime \prime}$ for some constants $C^{\prime}, C^{\prime \prime}>0$.

Furthermore, the phase has a unique critical point, given by $\left(u_{0}, \vartheta_{0}\right)=$ $\left(\nu / 2 \lambda\left(m_{x}\right), 0\right)$, and Hessian matrix of the phase $\Gamma_{\theta_{1}, \theta_{2}}(1, \cdot, \cdot)$ has determinant $-4 \lambda\left(m_{x}\right)^{2}$ and its signature is zero. Thus the critical point is nondegenerate, and the critical value is $\Gamma_{\theta_{1}, \theta_{2}}\left(1, u_{0}, 0\right)=0$. At the critical point, the exponent in the amplitude is

$$
\mathcal{E}_{1}\left(u_{0}, 0, \mathbf{v}_{1}, \mathbf{v}_{2}\right)=u_{0} \cdot \psi_{2}\left(\mathbf{v}_{1}, \mathbf{v}_{2}\right)
$$

When applying the Stationary Phase Lemma, at the $l$-th step we need to let the differential operator $k^{-l / 2} R_{\Gamma}^{l}$ act on the amplitude, where

$$
R_{\Gamma}:=\frac{i}{4 \cdot \lambda\left(m_{x}\right)} \cdot \frac{\partial^{2}}{\partial u \partial \vartheta},
$$

and evaluate the result at the critical point. One sees inductively that

$$
k^{-l / 2} R_{\Gamma}^{l}\left(e^{i \mathcal{E}_{1}\left(u, \vartheta, \mathbf{v}_{1}, \mathbf{v}_{2}\right)}\right)=H_{l} e^{i \mathcal{E}_{1}\left(u, \vartheta, \mathbf{v}_{1}, \mathbf{v}_{2}\right)}
$$

where $H_{l}$ is a polynomial of degree $\leq 3 l$ in $\left(\vartheta, \mathbf{v}_{1}, \mathbf{v}_{2}\right)$ and parity $(-1)^{l}$. The claimed asymptotic expansion follows by applying the Stationary Phase Lemma.

Lemma 4.3.5. As $k \rightarrow+\infty$, we have $J_{v_{1}, v_{2}}^{\prime \prime}(k)=O\left(k^{-\infty}\right)$.
Proof of Lemma 4.3.5. Let us view $J_{v_{1}, v_{2}}^{\prime \prime}(k)$ as an oscillatory integral in $\sqrt{k}$, with the phase $\Gamma_{\theta_{1}, \theta_{2}}(-1, u, \vartheta)$. We have

$$
\partial_{\vartheta} \Gamma_{\theta_{1}, \theta_{2}}(-1, u, \vartheta)=-2 u \cdot \lambda\left(m_{x}\right)-\nu \leq-\nu .
$$

The claim follows by integration by parts in $\vartheta$.

Lemma 4.3.6. $\mathcal{J}_{v_{1}, v_{2}}^{\prime \prime \prime}(k ; t, \delta)_{2}=O\left(k^{-1}\right)$ as $k \rightarrow+\infty$.
Proof of Lemma 4.3.6. Let us choose $\epsilon^{\prime} \in\left(0, \nu /\left(8 D \cdot \lambda\left(m_{x}\right)\right)\right)$, and consider the open cover $\mathcal{U}:=\left\{\left[-1,2 \epsilon^{\prime}\right),\left(\epsilon^{\prime}, 1\right]\right\}$. Let $\gamma_{1}(t)+\gamma_{2}(t)=1$ be a smooth partition of unity on $[-1,1]$ subordinate to $\mathcal{U}$. Thus

$$
\mathcal{J}_{v_{1}, v_{2}}^{\prime \prime \prime}(k)=\mathcal{J}_{v_{1}, v_{2}}^{\prime \prime \prime}(k)_{1}+\mathcal{J}_{v_{1}, v_{2}}^{\prime \prime \prime}(k)_{2},
$$

where $\mathcal{J}_{v_{1}, v_{2}}^{\prime \prime \prime}(k)_{j}$ is defined as $\mathcal{J}_{v_{1}, v_{2}}^{\prime \prime \prime}(k)$, except that the integrand has been further multiplied by $\gamma_{j}(t)$. Explicitly, let us write

$$
\mathcal{J}_{v_{1}, v_{2}}^{\prime \prime \prime}(k)_{j}=\int_{-\pi}^{\pi} \mathrm{d} \delta \int_{-1}^{1} \mathrm{~d} t\left[\mathcal{J}_{v_{1}, v_{2}}^{\prime \prime \prime}(k ; t, \delta)_{j}\right]
$$

where

$$
\begin{aligned}
\mathcal{J}_{v_{1}, v_{2}}^{\prime \prime \prime}(k ; t, \delta)_{j}= & \int_{0}^{+\infty} \mathrm{d} u \int_{-\infty}^{+\infty} \mathrm{d} \vartheta\left[e^{i \sqrt{k} \Gamma_{\theta_{1}, \theta_{2}}(t, u, \vartheta)} \cdot \gamma_{j}(t)\right. \\
& \left.\cdot \partial_{t}\left(e^{i u B_{t}+u\left[\psi_{2}\left(V_{1}, \mathbf{v}_{2}\right)-\frac{1}{2} A_{t}^{2}\right]}\right) \cdot \varrho_{1}(u) \cdot \varrho_{2}\left(k^{-\epsilon_{1}} \vartheta\right) \cdot u^{d-1}\right]
\end{aligned}
$$

Let us view the integral $\mathcal{J}_{v_{1}, v_{2}}^{\prime \prime \prime}(k ; t, \delta)_{j}$ as an oscillatory integral with phase $\Gamma_{\theta_{1}, \theta_{2}}(t, u, \vartheta)$.

On the support of $\gamma_{1}$, we have $u \leq 2 D$ and $t \leq \epsilon^{\prime}$; therefore,

$$
\partial_{\vartheta} \Gamma_{\theta_{1}, \theta_{2}}(t, u, \vartheta)=2 u \cdot t \cdot \lambda\left(m_{x}\right)-\nu \leq 4 D \cdot \epsilon^{\prime} \cdot \lambda\left(m_{x}\right)-\nu \leq-\frac{\nu}{2}
$$

Therefore, integration by parts in $\vartheta$ implies that $\mathcal{J}_{v_{1}, v_{2}}^{\prime \prime \prime}(k ; t, \delta)_{j}=O\left(k^{-\infty}\right)$, uniformly on the support of $\gamma_{1}$. It follows that $\mathcal{J}_{v_{1}, v_{2}}^{\prime \prime \prime}(k)_{1}=O\left(k^{-\infty}\right)$.

On the support of $\gamma_{2}$, on the other hand, $\Gamma_{\theta_{1}, \theta_{2}}(t, \cdot, \cdot)$ has the nondegenerate critical point

$$
(u(t), \vartheta(t))=\left(\frac{\nu}{2 t \lambda\left(m_{x}\right)}, 0\right)
$$

with Hessian matrix

$$
\operatorname{Hess}\left(\Gamma_{\theta_{1}, \theta_{2}}(t, \cdot, \cdot)\right)=\left(\begin{array}{cc}
0 & 2 t \cdot \lambda\left(m_{x}\right) \\
2 t \cdot \lambda\left(m_{x}\right) & 0
\end{array}\right)
$$

On the other hand, the integrand is divisible by $\vartheta$, hence it vanishes at the critical point. Furthermore, the integrand is of class $L^{1}$ as a function of the
parameter $t$, since the exponent getting differentiated is a smooth function of $t$ and $\sqrt{1-t^{2}}$.

Hence $\mathcal{J}_{v_{1}, v_{2}}^{\prime \prime \prime}(k, t, \delta)_{2}$ admits an asymptotic expansion in descending halfinteger powers of $k$, with leading power $k^{-1}$, and coefficients of class $L^{1}$ as functions of $t$, since the exponent getting differentiated is a smooth function of $t$ and $\sqrt{1-t^{2}}$.

Hence $\mathcal{J}_{v_{1}, v_{2}}^{\prime \prime \prime}(k ; t, \delta)_{2}$ admits an asymptotic expansion in descending halfinteger powers of $k$, with leading power $k^{-1}$, and coefficients of class $L^{1}$ as functions of $t$. The general term of the expansion will be a scalar multiple of

$$
\begin{equation*}
k^{-(l+1) / 2} \cdot R_{\Gamma_{t}}^{l}\left(\partial_{t}\left(\mathcal{E}_{t}\left(u, \vartheta, \mathbf{v}_{1}, \mathbf{v}_{2}\right)\right) \cdot e^{\mathcal{E}_{t}\left(u, \vartheta, \mathbf{v}_{1}, \mathbf{v}_{2}\right)}\right), \tag{4.23}
\end{equation*}
$$

where

$$
\mathcal{E}_{t}\left(u, \vartheta, \mathbf{v}_{1}, \mathbf{v}_{2}\right):=i u B_{t}+u\left[\psi_{2}\left(V_{t}, \mathbf{v}_{2}\right)-\frac{1}{2} A_{t}^{2}\right] .
$$

Given integers $a, b \geq 0$, let us denote by $H_{a, b}\left(\vartheta ; \mathbf{v}_{1} ; \mathbf{v}_{2}\right)$ a generic polynomial in $\left(\vartheta ; \mathbf{v}_{1}, \mathbf{v}_{2}\right)$, which is separately homogeneous of degree $a$ in $\vartheta$, and of degree $b$ in $\left(\mathbf{v}_{1}, \mathbf{v}_{2}\right)$, and by $H_{a}\left(\vartheta ; \mathbf{v}_{1}, \mathbf{v}_{2}\right)$ a generic polynomial in $\left(\vartheta ; \mathbf{v}_{1}, \mathbf{v}_{2}\right)$ homogeneous of degree $a$ (but perhaps not polyhomogeneous); both $H_{a, b}$ and $H_{a}$ are allowed to vary from line to line, and their coefficients depend smoothly on $u$. Thus $\mathcal{E}_{t}=u \cdot H_{2}=u \cdot\left(H_{2,0}+H_{1,1}+H_{0,2}\right), \partial_{t} \mathcal{E}_{t}=u \cdot\left(H_{2,0}+H_{1,1}\right)$ (here the polynomials do not depend on $u$ ). Hence we can split (4.23) as

$$
\begin{aligned}
& k^{-(l+1) / 2} \cdot\left[R_{\Gamma_{t}}^{l}\left(\rho(u) \cdot H_{2,0} \cdot e^{u \cdot\left(H_{2,0}+H_{1,1}+H_{0,2}\right)}\right)\right. \\
& \left.\quad+R_{\Gamma_{t}}^{l}\left(\rho(u) \cdot H_{1,1} \cdot e^{u \cdot\left(H_{2,0}+H_{1,1}+H_{0,2}\right)}\right)\right] .
\end{aligned}
$$

The proof of Lemma is completed by the following two claims, which can be proved inductively from the cases $l=0,1$ :

Claim 1. For $l=0,1,2, \ldots$

$$
R_{\Gamma_{t}}^{l}\left(\rho(u) \cdot H_{1,1} \cdot e^{\mathcal{E}_{t}}\right)
$$

is a sum of term of the form

$$
\left[H_{0,1} \cdot H_{p_{l}}+H_{1,1} \cdot H_{q_{l}}\right] \cdot e^{\mathcal{E}_{t}}
$$

where $p_{l}+1 \leq 3 l,(-1)^{p_{l}+1}=(-1)^{l}$, and $q_{l} \leq 3 l,(-1)^{q_{l}}=(-1)^{l}$.

Claim 2. For $l=0,1,2, \ldots$

$$
R_{\Gamma_{t}}^{l}\left(\rho(u) \cdot H_{2,0} \cdot e^{\mathcal{E}_{t}}\right)
$$

is a sum of term of the form $H_{a, b} \cdot e^{\mathcal{E}_{t}}$, where $b \leq 3 l,(-1)^{a+b}=(-1)^{l}$.

Since at the critical point $\vartheta=0$, the summands with a factor of the form $H_{a, b}$ with $a \geq 1$ all vanish at the critical point. It follows that the asymptotic expansion for (4.22) is as in the statement of Proposition 4.3.1. This completes the proof.

Let us next consider the asymptotics of $\Pi_{k \nu}\left(x_{1 k}, x_{2 k}\right)_{\pi}$. We shall prove the following analogue of the latter Proposition.
Proposition 4.3.2. Under the assumption of Theorem 1.2.3, as $k \rightarrow+\infty$ we have

$$
\begin{aligned}
\Pi_{k \nu}\left(x_{1 k}, x_{2 k}\right)_{\pi} \sim & \frac{1}{2 \lambda\left(m_{x}\right)} \cdot\left(\frac{\nu k}{2 \pi \lambda\left(m_{x}\right)}\right)^{d} \cdot e^{u_{0} \psi_{2}\left(\mathbf{v}_{2}, \mathbf{v}_{2}\right)} \\
& \cdot\left[1+\sum_{j \geq 1}^{+\infty} k^{-j / 2} A_{j}\left(x ; \mathbf{v}_{1}, \mathbf{v}_{2}\right)\right]
\end{aligned}
$$

where $A_{j}(x, \cdot, \cdot)$ is a polynomial of degree $\leq 3 j$ and parity $(-1)^{j}$.
Proof of Proposition 4.3.2. Since the proof is a slight modification of the one for Proposition 4.3.1, we shall be sketchy. In the integrand, $\varrho(\vartheta)$ has now been replaced by $\varrho_{\pi}(\vartheta)$, therefore integration in $\mathrm{d} \vartheta$ is restricted to $(\pi-2 \delta, \pi+2 \delta)$. Since $\varrho_{\pi}$ vanishes identically unless $-I_{2} \in G_{x}$, we shall assume that the latter condition holds. Let us set $\mathbf{v}_{1}^{\prime}:=\mathrm{d}_{m_{x}} \mu_{-I_{2}}\left(\mathbf{v}_{1}\right)$ and

$$
\begin{aligned}
x_{1 k} & :=\tilde{\mu}_{-I_{2}}\left(x_{1 k}\right) \\
& =x+\left(\frac{1}{\sqrt{k}} \theta_{1}+R_{3}\left(\frac{1}{\sqrt{k}} \theta_{1}\right), \frac{1}{\sqrt{k}} \mathbf{v}_{1}^{\prime}+R_{2}\left(\mathbf{v}_{1}\right)\right) .
\end{aligned}
$$

We may assume without loss that $\varrho_{\pi}(\vartheta)=\varrho_{0}(\vartheta-\pi)$. Let us perform the change of variable $\vartheta \mapsto \pi+\vartheta$, so that integration in $\mathrm{d} \vartheta$ is over $(-2 \delta, 2 \delta)$.

We have

$$
\begin{aligned}
& \Pi_{k \nu}\left(x_{1 k}, x_{2 k}\right)_{\pi} \\
& \sim \frac{k^{3 / 2} \nu}{2 \pi} \int_{0}^{+\infty} \mathrm{d} u \int_{-\infty}^{+\infty} \mathrm{d} \vartheta \int_{G / T} \mathrm{dV} \\
& G / T \\
& (g T)\left[\rho_{2}\left(k^{-\epsilon_{1}} \vartheta\right) \varrho_{1}(u) e^{i k \Gamma_{k}}\right. \\
& \left.\quad\left(e^{i(\pi+\vartheta / \sqrt{k})}-e^{-i(\pi+\vartheta / \sqrt{k})}\right) s\left(\tilde{\mu}_{g e^{-i \vartheta B / \sqrt{k}} g^{-1}}\left(x_{1 k}^{\prime}\right), x_{2 k}, k u\right)\right],
\end{aligned}
$$

where

$$
\begin{align*}
\Gamma_{k}\left(u, v_{1}, v_{2}, \vartheta, g T\right) & :=u \psi\left(\tilde{\mu}_{g e^{-i \vartheta B / \sqrt{k}} g^{-1}}\left(x_{1 k}^{\prime}\right), x_{2 k}\right)-\frac{\vartheta}{\sqrt{k}} \nu-\pi \nu  \tag{4.24}\\
& =\Psi_{k}\left(u, v_{1}^{\prime}, v_{2}, \vartheta, g T\right)-\pi \nu \tag{4.25}
\end{align*}
$$

Let us write $\Psi_{k}^{\prime}=\Psi_{k}\left(u, v_{1}^{\prime}, v_{2}, \vartheta, g T\right)$. Thus we may rewrite the latter integral in the following manner

$$
\begin{aligned}
& \Pi_{k \nu}\left(x_{1 k}, x_{2 k}\right)_{\pi} \\
& \sim e^{i \pi(1-k \nu)} \frac{k^{3 / 2} \nu}{2 \pi} \int_{0}^{+\infty} \mathrm{d} u \int_{-\infty}^{+\infty} \mathrm{d} \vartheta \int_{G / T} \mathrm{dV}_{G / T}(g T)\left[\rho_{2}\left(k^{-\epsilon_{1}} \vartheta\right) e^{i k \Psi_{k}^{\prime}}\right. \\
& \left.\quad \cdot \varrho_{1}(u)\left(e^{i(\pi+\vartheta / \sqrt{k})}-e^{-i(\pi+\vartheta / \sqrt{k})}\right) s\left(\tilde{\mu}_{g e^{-i \vartheta B / \sqrt{k}} g^{-1}}\left(x_{1 k}^{\prime}\right), x_{2 k}, k u\right)\right] \\
& \sim e^{i \pi(1-k \nu)} \cdot \Pi_{k \nu}\left(x_{1 k}^{\prime}, x_{2 k}\right)_{0}
\end{aligned}
$$

The statement of Proposition follows from Proposition 4.3.1.
Thus the proof is concluded.

### 4.4 Proof of Theorem 1.3.1

Before we delve into the proof let us recall some relevant concepts and results from [GS82c]. We shall use throughout the identification $T^{*} G \cong G \times \mathfrak{g}^{\vee}$ induced by right translations. If $R$ and $S$ are manifolds and $\Lambda \subset T^{*} R \times T^{*} S$ is a Lagrangian submanifold, the corresponding canonical relation is

$$
\Lambda^{\prime}:=\{((r, v),(s,-\gamma)):(r, v),(s, \gamma) \in \Lambda\}
$$

Definition 4.4.1. For every $f \in \mathcal{C}\left(\mathcal{O}_{\nu}\right)$, let $G_{f} \leq G$ be the stabilizer subgroup of $f$, and let $\mathfrak{g}_{f}$ be its Lie algebra. Let $H_{f}$ be the closed connected subgroup with Lie subalgebra $\mathfrak{h}_{f}:=\left\{\xi \in \mathfrak{g}_{f}:\langle f, \xi\rangle=0\right\}$. The locus

$$
\Lambda_{L}:=\left\{(g, r f) \in G \times \mathfrak{g}^{\vee}: f \in \mathcal{O}_{\nu}, r>0, g \in H_{f}\right\}
$$

is a Lagrangian submanifold of $T^{*} G$.
Definition 4.4.2. For every weight $\boldsymbol{\nu}$, Let us denote by $L=L_{\boldsymbol{\nu}}:=(k \boldsymbol{\nu})_{k=0}^{+\infty}$ the ladder sequence of weights generated by $\boldsymbol{\nu}$, and set

$$
\chi_{L}:=\sum_{k=1}^{+\infty} d_{k \nu} \chi_{k \nu} \in \mathcal{D}^{\prime}(G)
$$

Then we have the following
Theorem 4.4.1 (Theorem 6.3 of [GS82c]). $\chi_{L}$ is a Lagrangian distribution on $G$, and its associated conic Lagrangian submanifold of $T^{*} G \cong G \times \mathfrak{g}^{\vee}$ is $\Lambda_{L}$.

Consider the Hilbert space direct sum

$$
H(X)_{L}:=\bigoplus_{k=1}^{+\infty} H(X)_{k \nu}
$$

and let $\Pi_{L}: L^{2}(X) \rightarrow H(X)_{L}$ denote the corresponding orthogonal projector, $\Pi_{L}(\cdot, \cdot) \in \mathcal{D}^{\prime}(X \times X)$ its Schwartz kernel. then

$$
\Pi_{L}(x, y):=\int_{G} \overline{\chi_{L}(g)} \Pi\left(\tilde{\mu}_{g^{-1}}(x), y\right) \mathrm{d} V_{G}(g)
$$

We shall express $\Pi_{L}$ in functorial notation using functorial behaviour of wave fronts under pull-back and push-forward (see for instance $\S 1.3$ of [Dui96]) to draw conclusions on the singularities of $\Pi_{L}$.

To this end, let us consider the map

$$
f: G \times X \times X \rightarrow X \times X, \quad(g, x, y) \mapsto\left(\tilde{\mu}_{g^{-1}}(x), y\right)
$$

and the distribution $\hat{\Pi}:=f^{*}(\Pi) \in \mathcal{D}^{\prime}(G \times X \times X)$. Let

$$
\Sigma:=\{(x, r \alpha): x \in X, r>0\} \subset T^{*} X \backslash(0)
$$

denote the closed symplectic cone sprayed by the connection 1 -form; by [BdMS76], the wave front of $\Pi$ satisfies

$$
\mathrm{WF}^{\prime}(\Pi)=\operatorname{diag}(\Sigma) \subset \Sigma \times \Sigma .
$$

It follows that $\mathrm{WF}^{\prime}(\hat{\Pi}) \subseteq f^{*}(\operatorname{diag}(\Sigma))$. This implies the following.
Lemma 4.4.1. In terms of the identification $T^{*} G \cong G \times \mathfrak{g}^{\vee}$ induced by the right translations, the canonical relation of $\hat{\Pi}$ is

$$
\begin{aligned}
\mathrm{WF}^{\prime}(\hat{\Pi})=\{ & \left(\left(g, r \Phi_{G}\left(m_{x}\right)\right),\left(x, r \alpha_{x}\right),\left(y, r \alpha_{y}\right)\right) \in T^{*} G \times T^{*} X \times T^{*} X: \\
& \left.g \in G, x \in X, r>0, y=\tilde{\mu}_{g^{-1}}(x)\right\}
\end{aligned}
$$

where $m_{x}:=\pi(x)$.
Now let us give the functorial reformulation of the kernel of the ladder projector. Consider the diagonal map

$$
\Delta: G \times X \times X \rightarrow G \times G \times X \times X, \quad(g, x, y) \mapsto(g, g, x, y)
$$

and the projection

$$
p: G \times X \times X \rightarrow X \times X, \quad(g, x, y) \mapsto(x, y)
$$

Lemma 4.4.2. The Schwarts kernel $\Pi_{L} \in \mathcal{D}^{\prime}(X \times X)$ is given by

$$
\Pi_{L}=p_{*}\left(\Delta^{*}\left(\overline{\chi_{L}} \boxtimes \hat{\Pi}\right)\right) .
$$

Let $\sigma: T^{*} G \rightarrow T^{*} G$ be given by $(g, f) \mapsto(g,-f)$. Then

$$
\begin{aligned}
& \mathrm{WF}\left(\overline{\chi_{L}} \boxtimes \hat{\Pi}\right) \\
& \quad \subseteq\left(\sigma\left(\Lambda_{L}\right) \times(0)\right) \cup\left(\sigma\left(\Lambda_{L}\right) \times \mathrm{WF}(\hat{\Pi})\right) \cup((0) \times \mathrm{WF}(\hat{\Pi})) \\
& \quad \subset T^{*} G \times\left(T^{*} G \times T^{*} X \times T^{*} X\right)
\end{aligned}
$$

Therefore, the pull back $\Delta^{*}\left(\overline{\chi_{L}} \boxtimes \hat{\Pi}\right)$ is well-defined, and

$$
\begin{aligned}
& \mathrm{WF}\left(\Delta^{*}\left(\overline{\chi_{L}} \boxtimes \hat{\Pi}\right)\right) \subseteq \mathrm{d} \Delta^{*}\left(\overline{\chi_{L}} \boxtimes \hat{\Pi}\right) \\
& \quad \subseteq\left(\sigma\left(\Lambda_{L}\right) \times(0)\right) \cup \mathrm{d} \Delta^{*}\left(\sigma\left(\Lambda_{L}\right) \times \mathrm{WF}(\hat{\Pi})\right) \cup \mathrm{WF}(\hat{\Pi}) \\
& \quad \subset T^{*} G \times T^{*} X \times T^{*} X .
\end{aligned}
$$

Explicitly, we have

$$
\begin{aligned}
\mathrm{d} \Delta^{*}\left(\sigma\left(\Lambda_{L}\right) \times \mathrm{WF}(\hat{\Pi})\right)=\{ & \left(\left(g,-f+r \Phi_{G}\left(m_{x}\right)\right),\left(x, r \alpha_{x}\right),\left(y,-r \alpha_{y}\right)\right): \\
& \left.f \in \mathcal{C}(\mathcal{O}) g \in H_{f}, x \in X, r>0, y=\tilde{\mu}_{g^{-1}}(x)\right\}
\end{aligned}
$$

Using that $\Phi_{G}$ is nowhere vanishing, we can now apply Proposition 1.3.4 of [Dui96] to conclude the following.

Corollary 4.4.1. The wave front $\mathrm{WF}\left(\Pi_{L}\right) \subseteq\left(T^{*} X \backslash(0)\right) \times\left(T^{*} X \backslash(0)\right)$ of the distributional kernel $\Pi_{L}$ satisfies

$$
\mathrm{WF}\left(\Pi_{L}\right)=\left\{\left(\left(x, r \alpha_{x}\right),\left(y,-r \alpha_{y}\right)\right): \Phi_{G}\left(m_{x}\right) \in \mathcal{C}(\mathcal{O}), y \in H_{f} \cdot x\right\}
$$

where $H_{f} \cdot x$ is the $H_{f}$-orbit of $x$.
Corollary 4.4.2. Let $\mathrm{SS}\left(\Pi_{L}\right) \subseteq X \times X$ be the singular support of the distributional kernel $\Pi_{L}$. Then $\operatorname{SS}\left(\Pi_{L}\right) \subseteq \mathcal{Z}_{\nu}$.

### 4.4.1 The proof

Proof of Theorem 1.3.1. For every $\boldsymbol{\nu}=\left(\nu_{1}, \nu_{2}\right) \in \mathbb{Z}^{2}$ with $\nu_{1}>\nu_{2}$, let $P_{\boldsymbol{\nu}}$ : $L^{2}(X) \rightarrow L^{2}(X)_{\nu}$ be the orthogonal projector. Clearly

$$
\begin{equation*}
\Pi_{k \nu}=P_{k \nu} \circ \Pi_{L} \tag{4.26}
\end{equation*}
$$

In terms of the Schwartz kernels, (4.26) can be reformulated as follows:

$$
\begin{equation*}
\Pi_{k \nu}(x, y)=d_{k \nu} \int_{G} \mathrm{~d} V_{G}(g)\left[\overline{\chi_{k \nu}(g)} \Pi_{L}\left(\tilde{\mu}_{g^{-1}}(x), y\right)\right] \tag{4.27}
\end{equation*}
$$

Using the Weyl integration, character and dimension formulae, (4.27) can in turn be rewritten as follows (see chapter $\S 3$ ):

$$
\begin{equation*}
\Pi_{k \boldsymbol{\nu}}(x, y)=\frac{k\left(\nu_{1}-\nu_{2}\right)}{(2 \pi)^{2}} \int_{(-\pi, \pi)^{2}} \mathrm{~d} \boldsymbol{\vartheta}\left[e^{i k\langle\boldsymbol{\nu}, \boldsymbol{\vartheta}\rangle}\left(e^{i \vartheta \vartheta_{1}}-e^{i \vartheta_{2}}\right) F_{L}\left(x, y ; e^{i \boldsymbol{\vartheta}}\right)\right] \tag{4.28}
\end{equation*}
$$

where for $t \in T$ we set

$$
\begin{equation*}
F_{L}\left(x, y ; e^{i \vartheta}\right):=\int_{G / T} \mathrm{~d} V_{G / T}(g T)\left[\Pi_{L}\left(\tilde{\mu}_{g t^{-1} g^{-1}}(x), y\right)\right] \tag{4.29}
\end{equation*}
$$

Now suppose $K \Subset(X \times X) \backslash \mathcal{Z}_{\nu}$. We may assume without loss that $K$ is $G \times G$-invariant. There exist $G \times G$-invariant open subsets $A, B \subset X \times X$ such that

$$
K \subset A \Subset(X \times X) \backslash \mathcal{Z}_{\boldsymbol{\nu}}, \quad \mathcal{Z}_{\boldsymbol{\nu}} \subset B \Subset(X \times X) \backslash K, \quad X \times X=A \cup B
$$

Hence $A$ is a $G \times G$-invariant open neighbourhood of $K$ in $X \times X$, and the restriction of $\Pi_{L}$ to $A$ is $\mathcal{C}^{\infty}$.

Therefore, we get a $\mathcal{C}^{\infty}$ function

$$
R: T \times G / T \times A \rightarrow \mathbb{C}, \quad(t, g T,(x, y)) \mapsto \Pi_{L}\left(\tilde{\mu}_{g t^{-1} g^{-1}}(x), y\right)
$$

With $F_{L}$ as in (4.29), we obtain a $\mathcal{C}^{\infty}$ function on $T \times A$ by setting

$$
\beta:(t,(x, y)) \mapsto \Delta(t) F_{L}(x, y ; t) .
$$

Let us denote with $\mathcal{F}_{T}$ the Fourier transform with respect to $t \in T$ of a function on $T \times A$, viewed as a function on $\mathbb{Z}^{2} \times A$; then (4.28) may be rewritten

$$
\begin{equation*}
\Pi_{k \nu}(x, y)=\frac{k}{2}\left(\nu_{1}-\nu_{2}\right) \cdot \mathcal{F}_{T}(\beta)(k \boldsymbol{\nu} ; x, y) \tag{4.30}
\end{equation*}
$$

The statement of Theorem 1.3.1 follows from (4.30) and the previous considerations.

### 4.5 Proof of Theorem 1.3.2

We shall assume in this section that the assumptions of Theorem 1.3.2 hold.

### 4.5.1 Preliminaries

Before attacking the proof, it is in order to list some useful preliminaries (see also the discussion of $\S 2$ of [Pao12]).

For any $m \in M$, let $\operatorname{val}_{m}: \mathfrak{g} \rightarrow T_{m} M$ be the evaluation map $\xi \mapsto \xi_{M}(m)$, and similarly for any $x \in X$ let $\operatorname{val}_{x}: \mathfrak{g} \rightarrow T_{x} X$ be the evaluation map $\xi \mapsto \xi_{X}(x)$.

## Ray transversality and locally free actions

Since $\tilde{\mu}$ preserves the connection 1-form, the induced action of $G$ on $T^{*} X$ leaves the symplectic cone $\Sigma$ invariant. The restricted action is of course Hamiltonian, and its moment map $\tilde{\Phi}_{G}: \Sigma \rightarrow \mathfrak{g}$ is the restriction to $\Sigma$ of the cotangent Hamiltonian map on $T^{*} X$.

If $m \in M_{\mathcal{O}_{\nu}}^{G}$, then by equivariance $\Phi_{G}$ is transverse to $\mathbb{R}_{+} \cdot \Phi_{G}(m)$. Hence,

$$
\begin{equation*}
\mathrm{d}_{m} \Phi_{G}\left(T_{m} M\right)+\operatorname{span}\left(\Phi_{G}(m)\right)=\mathfrak{g} . \tag{4.31}
\end{equation*}
$$

Suppose $x \in \pi^{-1}(m) \subset X$ and $r>0$, and consider $\sigma=\left(x, r \alpha_{x}\right) \in \Sigma$. Then it follows from (4.31) that

$$
\mathrm{d}_{\sigma} \tilde{\Phi}_{G}\left(T_{\sigma} \Sigma\right)=\mathrm{d}_{m} \Phi_{G}\left(T_{m} M\right)+\operatorname{span}\left(\Phi_{G}(m)\right)=\mathfrak{g} .
$$

Thus $\tilde{\Phi}_{G}$ is submersive at any $\left(x, r \alpha_{x}\right)$ with $x \in X_{\mathcal{O}_{\nu}}^{G}$. If we let $\Sigma_{\mathcal{O}_{\nu}}^{G} \cong$ $X_{\mathcal{O}_{\nu}}^{G} \times \mathbb{R}_{+}$denote the inverse image of $X_{\mathcal{O}_{\nu}}^{G}$ in $\Sigma$, we conclude therefore that $G$ acts locally free on $\Sigma_{\mathcal{O}_{\nu}}^{G}$, and this clearly implies that it acts locally freely on $X_{\mathcal{O}_{\nu}}^{G}$.

The previous implications may obviously be reversed, and we obtain the following.

Lemma 4.5.1. The following conditions are equivalent

1. $\Phi_{G}$ is transverse to $\mathbb{R}_{+} \cdot i \boldsymbol{\nu}$;
2. $\tilde{\mu}$ is locally free on $X_{O_{\nu}}^{G}$;
3. for any $x \in X_{\mathcal{O}_{\nu}}^{G}, \operatorname{val}_{x}$ is injective;
4. for any $m \in M_{\mathcal{O}_{\nu}}^{G}, \operatorname{val}_{m}$ is injective on $\Phi_{G}(m)^{\perp_{\mathfrak{g}}}$.

## The vector field $\Upsilon_{\mu, \nu}$

Let us consider the normal vector field $\Upsilon=\Upsilon_{\mu, \nu}$ to $M_{\mathcal{O}_{\nu}}^{G}$ appearing in the statement of Theorem 1.3.2.

By definition, $m \in M_{\mathcal{O}_{\nu}}^{G}$ if and only if $\Phi_{G}(m)$ is similar to $i \lambda_{\boldsymbol{\nu}}(m) D_{\boldsymbol{\nu}}$, for some $\lambda_{\nu}(m)>0$. Equating norms and traces, we obtain

$$
\lambda_{\boldsymbol{\nu}}(m)=\frac{\left\|\Phi_{G}(m)\right\|}{\|\boldsymbol{\nu}\|}=-i \frac{\operatorname{trace}\left(\Phi_{G}(m)\right)}{\nu_{1}+\nu_{2}} \quad\left(m \in M_{\mathcal{O}_{\boldsymbol{\nu}}}^{G}\right) .
$$

Since $\nu_{1}>\nu_{2}$, there exists a unique coset $h_{m} T \in G / T$ such that

$$
\begin{equation*}
\Phi_{G}(m)=i \lambda_{\nu}(m) h_{m} D_{\nu} h_{m}^{-1} \tag{4.32}
\end{equation*}
$$

Let us set $\boldsymbol{\nu}_{\perp}:=\left(-\nu_{2}, \nu_{1}\right)^{t}$, and define $\boldsymbol{\rho}=\boldsymbol{\rho}_{\nu}: M_{\mathcal{O}_{\nu}}^{G} \rightarrow \mathfrak{g}$ by setting

$$
\begin{equation*}
\boldsymbol{\rho}(m):=i h_{m} D_{\boldsymbol{\nu}_{\perp}} h_{m}^{-1} \quad\left(m \in M_{\mathcal{O}_{\nu}}^{G}\right) . \tag{4.33}
\end{equation*}
$$

Then $\boldsymbol{\rho}(m)_{M} \in \mathfrak{X}(M)$ is the vector field on $M$ induced by $\boldsymbol{\rho}(m) \in \mathfrak{g}$; its evaluation at $m^{\prime} \in M$ is $\boldsymbol{\rho}(m)_{M}\left(m^{\prime}\right)$ (and similarly for $X$ ).

Definition 4.5.1. The vector field $\Upsilon=\Upsilon_{\mu, \nu}$ along $M_{\mathcal{O}_{\nu}}^{G}$ is

$$
\Upsilon(m):=J_{m}\left(\boldsymbol{\rho}(m)_{M}(m)\right) \quad\left(m \in M_{\mathcal{O}_{\nu}}^{G}\right)
$$

With abuse of notation, recalling (1.3) we shall also denote by $\Upsilon$ the vector field along $X_{\mathcal{O}_{\nu}}^{G}$ given by

$$
\Upsilon(x):=J_{x}\left(\boldsymbol{\rho}\left(m_{x}\right)_{X}(x)\right), \quad m_{x}:=\pi(x) .
$$

Notice that

$$
\left\langle\Phi_{G}(m), \boldsymbol{\rho}(m)\right\rangle=\lambda_{\boldsymbol{\nu}}(m)\left\langle\boldsymbol{\nu}, \boldsymbol{\nu}_{\perp}\right\rangle=0 \quad\left(m \in M_{\mathcal{O}_{\boldsymbol{\nu}}}^{G}\right)
$$

Therefore, in view of (1.4) for any $x \in \pi^{-1}(m)$ we have

$$
\boldsymbol{\rho}(m)_{X}(x)=\boldsymbol{\rho}(m)_{M}^{\sharp}(x) .
$$

## A spectral characterization of $G \cdot M_{\nu}^{G}$

Suppose that $-i \Phi_{G}(m)$ has eigenvalues $\lambda_{1}(m) \geq \lambda_{2}(m)$. Then $m \in M_{\mathcal{O}_{\nu}}^{G}$ if and only if $\lambda_{1}(m) \nu_{2}-\lambda_{2}(m) \nu_{1}=0$. We shall give a similar spectral characterization of $G \cdot M_{\nu}^{T}$. Notice that if $\lambda_{1}(m)=\lambda_{2}(m)$, then $\Phi_{G}(m)$ is a multiple of the identity, hence certainly $m \notin M_{\mathcal{O}_{\nu}}^{G}$. Thus we may as well assume that $\lambda_{1}(m)>\lambda_{2}(m)$.

Proposition 4.5.1. Suppose $m \in M$, and let the eigenvalues of $-i \Phi_{G}(m)$ be $\lambda_{1}(m)>\lambda_{2}(m)$. Then $m \in G \cdot M_{\nu}^{T}$ if and only if

$$
\begin{equation*}
t(m, \boldsymbol{\nu}):=\frac{\lambda_{1}(m) \nu_{2}-\lambda_{2}(m) \nu_{1}}{\left(\nu_{1}+\nu_{2}\right)\left(\lambda_{1}(m)-\lambda_{2}(m)\right)} \in[0,1 / 2) \tag{4.34}
\end{equation*}
$$

Proof of Proposition 4.5.1. Let us set $\boldsymbol{\lambda}(m):=\left(\lambda_{1}(m), \lambda_{2}(m)\right)$, and let $D_{\boldsymbol{\lambda}}$ be the corresponding diagonal matrix. By definition, $m \in G \cdot M_{\nu}^{T}$ if and only if there exists $g \in S U(2) \leq G$ such that $\operatorname{diag}\left(g D_{\boldsymbol{\lambda}} g^{-1}\right) \in \mathbb{R}_{+} \cdot \boldsymbol{\nu}$. This is equivalent to the condition that there exists $u, v \in \mathbb{C}$ such that

$$
\left(\begin{array}{cc}
u & -\bar{w}  \tag{4.35}\\
w & \bar{u}
\end{array}\right) D_{\lambda}\left(\begin{array}{cc}
\bar{u} & \bar{w} \\
-w & u
\end{array}\right)=c\left(\begin{array}{cc}
\nu_{1} & a \\
\bar{a} & \nu_{2}
\end{array}\right)
$$

for some $c>0$ and $a \in \mathbb{C}$. If we set $t=|w|^{2}$, we conclude that $m \in G \cdot M_{\nu}^{T}$ if and only if there exists $t \in[0,1]$ such that

$$
\boldsymbol{\lambda}_{t}(m):=\binom{(1-t) \lambda_{1}(m)+t \lambda_{2}(m)}{(1-t) \lambda_{2}(m)+t \lambda_{1}(m)} \in \mathbb{R}_{+} \cdot\binom{\nu_{1}}{\nu_{2}}
$$

The condition $\boldsymbol{\lambda}_{t}(m) \wedge \boldsymbol{\nu}=0$ translates into the equality $t=t(m, \boldsymbol{\nu})$. Hence we need to have $t(m, \boldsymbol{\nu}) \in[0,1]$. Given this, $\boldsymbol{\lambda}_{t}(m)$ is a positive multiple of $\boldsymbol{\nu}$ if and only if
$(1-t(m, \boldsymbol{\nu})) \lambda_{1}(m)+t(m, \boldsymbol{\nu}) \lambda_{2}(m)>t(m, \boldsymbol{\nu}) \lambda_{1}(m)+(1-t(m, \boldsymbol{\nu})) \lambda_{2}(m)$,
and this is equivalent to $t(m, \boldsymbol{\nu})<1 / 2$.
Conversely, suppose that $t(m, \boldsymbol{\nu}) \in[0,1 / 2)$, and define

$$
g:=\left(\begin{array}{cc}
\sqrt{1-t(m, \boldsymbol{\nu})} & -\sqrt{t(m, \boldsymbol{\nu})} \\
\sqrt{t(m, \boldsymbol{\nu})} & \sqrt{1-t(m, \boldsymbol{\nu})}
\end{array}\right) .
$$

### 4.5.2 The proof

Proof of Theorem 1.3.2. As $\Phi_{G}$ is equivariant, it is transverse to $\mathbb{R}_{+} \cdot i D_{\nu}$ if and only if it is transverse to $\mathbb{R}_{+} \cdot \mathcal{O}_{\nu}$. Given that $\nu_{1}>\nu_{2}, \mathcal{O}_{\nu}$ is 2dimensional (and diffeomorphic to $S^{2}$ ); therefore, $\mathbb{R}_{+} \cdot \mathcal{O}_{\nu}$ has codimension 1 in $\mathfrak{g}$. Similarly, $\mathbb{R}_{+} \cdot i D_{\nu}$ has codimension 1 in $\mathfrak{t}^{\vee}$. Given $\mathbf{0} \notin \Phi_{G}(M)$, we conclude the following.
Step 4.5.1. $M_{\nu}^{G}, M_{\mathcal{O}_{\nu}}^{G}$ and $M_{\nu}^{T}$ are compact and smooth (real) submanifolds of $M . M_{\nu}^{G}$ has codimension 3 , and $M_{\mathcal{O}_{\nu}}^{G}$ and $M_{\nu}^{T}$ are hypersurfaces.

The Weyl chambers in $\mathfrak{t}$ are the half-planes

$$
\mathfrak{t}_{+}:=\left\{\boldsymbol{\mu}: \mu_{1}>\mu_{2}\right\}, \quad \mathfrak{t}_{-}:=\left\{\boldsymbol{\mu}: \mu_{1}<\mu_{2}\right\}
$$

and clearly with our identifications $i D_{\nu} \leftrightarrow \boldsymbol{\nu} \in \mathfrak{t}_{+}$. Since $\Phi_{G}(m) \cap \mathfrak{t}_{+}$is a convex polytope ([GS82a], [GS84], [Kir84a]), $\Phi_{G}(M) \cap \mathbb{R}_{+} \cdot i D_{\nu}$ is a closed segment $I$. Furthermore, for any $a \in I$, the inverse image $\Phi_{G}^{-1}(a) \subseteq M$ is also connected ([Kir84b], [Ler95]). Thus we obtain the following conclusion.

Step 4.5.2. $M_{\nu}^{G}, M_{\mathcal{O}_{\nu}}^{G}$ and $M_{\nu}^{T}$ are connected.
Proof of Step 4.5.2. The previous considerations imply easily that $M_{\nu}^{G}$ is connected. Given this, since $M_{\mathcal{O}_{\nu}}^{G}=G \cdot M_{\nu}^{G}$ the connectedness of $G$ implies the one of $M_{\mathcal{O}_{\nu}}^{G}$. Let us consider $M_{\nu}^{T}$. Since $\Phi_{T}(M)$ is a convex polytope ([GS82a], [Ati82]), $\Phi_{T}(M) \cap \mathbb{R}_{+} \cdot i D_{\nu}$ is also a connected segment $I^{\prime}$. The statement follows since the fibers of $\Phi_{T}$ are connected again by [Kir84b] and [Ler95].

For any $m \in M_{\mathcal{O}_{\nu}}^{G}$, let us set

$$
M_{\Phi_{G}(m)}^{G}:=\Phi_{G}^{-1}\left(\mathbb{R}_{+} \cdot \Phi_{G}(m)\right) .
$$

Since $\Phi_{G}$ is transverse to $\mathbb{R}_{+} \cdot \boldsymbol{\nu}$, by equivariance it is also transverse to $\mathbb{R}_{+} \cdot \Phi_{G}(m)$; hence $M_{\Phi_{G}(m)}^{G}$ is also a connected real submanifold of $M$, of real codimension 3 and contained in $M_{\mathcal{O}_{\nu}}^{G}$.

Let us consider the normal bundle $N\left(M_{\Phi_{G}(m)}^{G}\right)$ to $M_{\Phi_{G}(m)}^{G} \subset M$. For any $\xi \in \mathfrak{g}$, let $\xi^{\perp} \subset \mathfrak{g}$ be the orthocomplement to $\xi$. Under the equivariant identification $\mathfrak{g} \cong \mathfrak{g}^{\vee}, \xi^{\perp}$ corresponds to $\xi^{0}$.

For any subset $L \subseteq \mathfrak{g}$, let $L^{\perp_{\mathfrak{g}}}$ denote the orthocomplement of $L$ (that is, of the linear span of $L$ ) under the pairing $\langle\cdot, \cdot\rangle_{\mathfrak{g}}$.
Lemma 4.5.2. For any $m \in M_{O_{\nu}}^{G}$, we have

$$
N_{m}\left(M_{\Phi_{G}(m)}^{G}\right)=J_{m} \circ \operatorname{val}_{m}\left(\Phi_{G}(m)^{\perp_{\mathfrak{g}}}\right) .
$$

Similarly, for any $m \in M_{\nu}^{T}$, we have

$$
N_{m}\left(M_{\boldsymbol{\nu}}^{T}\right)=J_{m} \circ \operatorname{val}_{m}\left((i \boldsymbol{\nu})^{\perp_{\mathrm{t}}}\right)=J_{m} \circ \operatorname{val}_{m}\left(i \boldsymbol{\nu}_{\perp}\right)
$$

Proof of Lemma 4.5.2. If $v \in T_{m} M_{\Phi_{G}(m)}^{G}$, then $\mathrm{d}_{m} \Phi_{G}(v)=a \Phi_{G}(m)$ for some $a \in \mathbb{R}$. Given $\eta \in \Phi_{G}(m)^{\perp_{\mathfrak{g}}}$, and with $\rho$ as in (1.1), we have

$$
\begin{aligned}
\rho_{m}\left(J_{m}\left(\eta_{M}(m)\right), v\right) & =\omega_{m}\left(\eta_{M}(m), v\right)=\mathrm{d}_{m} \Phi^{\eta}(v) \\
& =\left\langle\mathrm{d}_{m} \Phi(v), \eta\right\rangle_{\mathfrak{g}}=a\left\langle\Phi_{G}(m), \eta\right\rangle_{\mathfrak{g}}=0 .
\end{aligned}
$$

Therefore $J_{m} \circ \operatorname{val}_{m}\left(\Phi_{G}(m)^{\perp_{\mathfrak{g}}}\right) \subseteq N_{m}\left(M_{\Phi_{G}(m)}^{G}\right)$. Since we have that both $\Phi_{G}^{\perp_{\mathfrak{g}}}$ and $N_{m}\left(M_{\Phi_{G}(m)}^{G}\right)$ are 3-dimensional, it suffices to recall that by Lemma 4.5.1 val $_{m}$ is injective when restricted to $\Phi_{G}(m)^{\perp_{\mathfrak{g}}}$.

The proof of the second statement is similar.
For any vector subspace $L \subseteq \mathfrak{g}$, let us set $L_{M}(m):=\operatorname{val}_{m}(L) \subseteq T_{m} M$ $(m \in M)$. For any $m \in M_{\mathcal{O}_{\nu}}^{G}$, given that $M_{\mathcal{O}_{\nu}}^{G}$ is the $G$-saturation of $M_{\Phi_{G}(m)}^{G}$, we have

$$
\begin{equation*}
T_{m} M_{\mathcal{O}_{\nu}}^{G}=T_{m} M_{\Phi_{G}(m)}^{G}+\mathfrak{g}_{M}(m) \tag{4.36}
\end{equation*}
$$

Therefore, passing to $\rho_{m}$-orthocomplements

$$
N_{m}\left(M_{\mathcal{O}_{\nu}}^{G}\right)=N_{m}\left(M_{\Phi_{G}(m)}^{G}\right) \cap \mathfrak{g}_{M}(m)^{\perp_{\rho_{m}}} .
$$

We conclude from relation (4.36) and Lemma 4.5.2 that the normal bundle $N_{m}\left(M_{\mathcal{O}_{\nu}}^{G}\right)$ is the set of all vectors $J_{m}\left(\eta_{M}(m)\right) \in T_{m} M$, with $\eta \in \Phi_{G}(m)^{\perp_{\mathfrak{g}}}$, such that $\rho_{m}\left(J_{m}\left(\eta_{M}(m)\right), \xi_{M}\right)=0$ for every $\xi \in \mathfrak{g}$. From this remark we can draw the following conclusion.
Step 4.5.3. Let $\Upsilon=\Upsilon_{\mu, \nu}$ be as in $\S 4.5 .1$. Then for any $m \in M_{\mathcal{O}_{\nu}}^{G}$ we have

$$
N_{m}\left(M_{\mathcal{O}_{\nu}}^{G}\right)=\operatorname{span}(\Upsilon(m)) .
$$

In particular, $M_{O_{\nu}}^{G}$ is orientable.
Proof of Step 4.5.3. By the above,

$$
\begin{aligned}
& N_{m}\left(M_{\mathcal{O}_{\nu}}^{G}\right) \\
& =\left\{J_{m}\left(\eta_{M}(m)\right) ; \eta \in \Phi_{G}(m)^{\perp_{\mathfrak{g}}} \wedge \rho_{m}\left(J_{m}\left(\eta_{M}(m)\right), \xi_{M}(m)\right)=0, \forall \xi \in \mathfrak{g}\right\} \\
& =\left\{J_{m}\left(\eta_{M}(m)\right) ; \eta \in \Phi_{G}(m)^{\perp_{\mathfrak{g}}} \wedge \omega_{m}\left(\eta_{M}(m), \xi_{M}(m)\right)=0, \forall \xi \in \mathfrak{g}\right\} \\
& =\left\{J_{m}\left(\eta_{M}(m)\right) ; \eta \in \Phi_{G}(m)^{\perp_{\mathfrak{g}}} \wedge \eta_{M}(m) \in \operatorname{ker}\left(\mathrm{d}_{m} \Phi_{G}\right), \forall \xi \in \mathfrak{g}\right\} \\
& =\left\{J_{m}\left(\eta_{M}(m)\right) ; \eta \in \Phi_{G}(m)^{\perp_{\mathfrak{g}}} \wedge\left[\eta, \Phi_{G}(m)\right]=0, \forall \xi \in \mathfrak{g}\right\} .
\end{aligned}
$$

The latter equality holds because, by the equivariance of $\Phi_{G}$, we have

$$
\begin{aligned}
\mathrm{d}_{m} \Phi_{G}\left(\eta_{M}(m)\right) & =\left.\frac{\mathrm{d}}{\mathrm{~d} t} \Phi_{G}\left(\mu_{e^{t \eta}}(m)\right)\right|_{t=0}=\left.\frac{\mathrm{d}}{\mathrm{~d} t} \operatorname{Ad}_{e^{t \eta}} \Phi_{G}(m)\right|_{t=0} \\
& =\left[\eta, \Phi_{G}(m)\right]
\end{aligned}
$$

There exists a unique $h_{m} T \in G / T$ such that $\Phi_{G}(m)=i \lambda_{\nu}(m) h_{m} D_{\nu} h_{m}^{-1}$. It is then clear that $\left\langle\Phi_{G}(m), \eta\right\rangle_{\mathfrak{g}}=0$ and $\left[\eta, \Phi_{G}(m)\right]=0$ if and only if

$$
\eta \in \operatorname{span}\left(i h_{m} D_{\boldsymbol{\nu}_{\perp}} h_{m}^{-1}\right)=\operatorname{span}(\boldsymbol{\rho}(m))
$$

where $\boldsymbol{\rho}(m)$ is as in (4.33). This completes the proof on Step 4.5.3.
Step 4.5.4. $M_{\mathcal{O}_{\nu}}^{G} \cap M_{\nu}^{T}=M_{\nu}^{G}$.
Proof of Step 4.5.4. Obviously, $M_{\mathcal{O}_{\nu}}^{G} \cap M_{\nu}^{T} \supseteq M_{\nu}^{G}$. Conversely, suppose $m \in M_{\mathcal{O}_{\nu}}^{G} \cap M_{\nu}^{T}$. Then on the one hand $\Phi_{G}(m)$ is similar to a positive multiple of $i D_{\nu}$; for a unique $h_{m} T \in G / T$,

$$
\Phi_{G}(m)=i \lambda_{\nu}(m) h_{m} D_{\nu} h_{m}^{-1}
$$

where we can assume without loss that $h_{m} \in S U(2)$. On the other hand $\operatorname{diag}\left(\Phi_{G}(m)\right)$ is a positive multiple of $i \boldsymbol{\nu}$. Hence the diagonal of $h_{m} D_{\boldsymbol{\nu}} h_{m}^{-1}$ is a positive multiple of $\boldsymbol{\nu}$. Let us write $h_{m}$ as in (4.35), and argue as in the proof of Proposition 4.5.1; using that $\nu_{1}^{2} \neq \nu_{2}^{2}$, one concludes readily that $h_{m}$ is diagonal. Hence $h_{m} D_{\nu} h_{m}^{-1}=D_{\boldsymbol{\nu}}$, and so $\Phi_{G}(m) \in \mathbb{R}_{+} \cdot i \boldsymbol{\nu}$. Thus $m \in M_{\nu}^{G}$.
Step 4.5.5. For any $m \in M_{\nu}^{G}, T_{m} M_{\mathcal{O}_{\nu}}^{G}=T_{m} M_{\nu}^{T}$.
Proof of Step 4.5.5. If $m \in M_{\nu}^{G}$, then $h_{m}$ is the identity matrix in (4.32) and (4.33); therefore, $\Upsilon(m)=J_{m}\left(\left(i D_{\nu_{\perp}}\right)_{M}(m)\right)$. Hence, we obtain that $N_{m}\left(M_{\mathcal{O}_{\nu}}^{G}\right)=\operatorname{span}\left(J_{m}\left(\left(i D_{\nu_{\perp}}\right)(m)\right)\right)$. The claim follows from this remark and Lemma 4.5.2.

Step 4.5.6. For any $M_{\mathcal{O}_{\nu}}^{G}=\partial\left(G \cdot M_{\nu}^{T}\right)$.
Proof of Step 4.5.6. Suppose $m \in M_{\mathcal{O}_{\nu}}^{G}$. Thus $\Phi_{G}(m)=i \lambda(m) h_{m} D_{\nu} h_{m}^{-1}$ for a unique $h_{m} T \in G / T$. Let us choose $\delta>0$ arbitrarily small, and let $M(m, \delta) \subseteq M$ be the open ball centred at $m$ and radius $\delta$ in the Riemannian distance on $M$. Since $\Phi_{G}$ is transverse to $\mathbb{R}_{+} \cdot i \boldsymbol{\nu}$, there exists $\epsilon_{1}>0$ such
that the following holds. For every $\epsilon \in\left(-\epsilon_{1}, \epsilon_{1}\right)$ there exists $m^{\prime} \in M(m, \delta)$ with

$$
\begin{equation*}
\Phi_{G}\left(m^{\prime}\right)=i \lambda\left(m^{\prime}\right) h_{m} D_{\boldsymbol{\nu}+\epsilon \boldsymbol{\nu}_{\perp}} h_{m}^{-1} \tag{4.37}
\end{equation*}
$$

for some $\lambda\left(m^{\prime}\right)>0$ (see $\S 2$ of [Pao12]). This implies that the eigenvalues of $-i \Phi_{G}\left(m^{\prime}\right)$ are

$$
\lambda_{1}\left(m^{\prime}\right):=\lambda\left(m^{\prime}\right)\left(\nu_{1}-\epsilon \nu_{2}\right), \quad \lambda_{2}\left(m^{\prime}\right):=\lambda\left(m^{\prime}\right)\left(\nu_{2}+\epsilon \nu_{1}\right) .
$$

Therefore, the invariant defined in (4.34) takes the following value at $m^{\prime}$ :

$$
t\left(m^{\prime}, \boldsymbol{\nu}\right)=-\frac{\epsilon}{\nu_{1}+\nu_{2}} \cdot \frac{\nu_{1}^{2}+\nu_{2}^{2}}{\left(\nu_{1}-\nu_{2}\right)-\epsilon\left(\nu_{1}+\nu_{2}\right)} .
$$

Therefore, if $\epsilon\left(\nu_{1}+\nu_{2}\right)>0$ (and $\epsilon$ is sufficiently small) then $m^{\prime} \notin G \cdot M_{\nu}^{T}$ by Proposition 4.5.1. This implies $M_{\mathcal{O}_{\nu}}^{G} \subseteq \partial\left(G \cdot M_{\nu}^{T}\right)$.

To prove the reverse inclusion, assume that $m \in G \cdot M_{\nu}^{T} \backslash M_{\mathcal{O}_{\nu}}^{G}$. Then $t(m, \boldsymbol{\nu}) \in[0,1 / 2)$ by 4.5.1. Furthermore, $t(m, \boldsymbol{\nu}) \neq 0$, since otherwise $m \in M_{\mathcal{O}_{\nu}}^{G}$. Hence $t(m, \boldsymbol{\nu}) \in(0,1 / 2) ;$ by continuity, then, $t\left(m^{\prime}, \boldsymbol{\nu}\right) \in(0,1 / 2)$ for every $m^{\prime}$ in a sufficiently small open neighbourhood of $m$. Hence Proposition 4.5.1 implies that $G \cdot M_{\nu}^{T} \backslash M_{\mathcal{O}_{\nu}}^{G}$ contains an open neighborhood of $m$ in $M$. Thus $G \cdot M_{\nu}^{T} \backslash M_{\mathcal{O}_{\nu}}^{G}$ is open, and in particular $m \notin \partial\left(G \cdot M_{\nu}^{T}\right)$.

Step 4.5.7. $\Upsilon$ is outer oriented if $\nu_{1}+\nu_{2}>0$ and inner oriented if $\nu_{1}+\nu_{2}<0$.
Proof of Step 4.5.7. Let denote by $\mathcal{B}_{\nu}$ the collection of all $B \in \mathfrak{g}$ such that $\operatorname{diag}\left(g B g^{-1}\right) \in \mathbb{R}_{+} \cdot i \boldsymbol{\nu}$ for some $g \in G$. Thus $\mathcal{B}_{\boldsymbol{\nu}}$ is a conic and invariant closed subset of $\mathfrak{g} \backslash\{0\}$; in addition, $m \in G \cdot M_{\nu}^{T}$ if and only if $\Phi_{G}(m) \in \mathcal{B}_{\nu}$.

If $\lambda_{1}(B) \geq \lambda_{2}(B)$ are the eigenvalues of $-i B$, the Proposition 4.5.1 implies that $B \in \mathcal{B}_{\nu}$ if and only if $\lambda_{1}(B)>\lambda_{2}(B)$ and

$$
t(B, \boldsymbol{\nu}):=\frac{\lambda_{1}(B) \nu_{2}-\lambda_{2}(B) \nu_{1}}{\left(\nu_{1}+\nu_{2}\right)\left(\lambda_{1}(B)-\lambda_{2}(B)\right)} \in[0,1 / 2)
$$

In particular, if $t(B, \boldsymbol{\nu}) \in(0,1 / 2)$ then $B$ belongs to the interior of $\mathcal{B}_{\boldsymbol{\nu}}$.
Suppose $m \in M_{\nu}^{G}$ and consider the path

$$
\gamma_{1}:(-\epsilon, \epsilon) \ni \tau \mapsto \Phi_{G}(m+\tau \Upsilon(m)) \in \mathfrak{g}
$$

defined for sufficiently small $\epsilon>0$; the expression $m+\tau \Upsilon(m) \in M$ is meant in an adapted coordinate system on $M$ centred at $m$. Then

$$
\begin{aligned}
& \gamma_{1}(0)=\Phi_{G}(m)=i \lambda_{\boldsymbol{\nu}}(m) D_{\boldsymbol{\nu}} \\
& \dot{\gamma}_{1}(0)=\omega_{m}(\cdot, \Upsilon(m))=\rho_{m}\left(\cdot,\left(i D_{\boldsymbol{\nu}_{\perp}}\right)_{M}(m)\right)
\end{aligned}
$$

Let us consider a smooth positive function $y:(-\epsilon, \epsilon) \rightarrow \mathbb{R}$, to be determined but subject to the condition $y(0)=\lambda_{\nu}(m)$. Let us define a second path of the form

$$
\gamma_{2}(\tau):=i y(\tau) \operatorname{Ad}_{e^{\tau \xi}}\left(i D_{\nu+a \tau \nu_{\perp}}\right)
$$

where $a>0$ is a constant also to be determined.
Then

$$
\begin{aligned}
& \gamma_{1}(0)=\gamma_{2}(0) \\
& \dot{\gamma}_{2}(0)=i\left[\dot{y}(0) D_{\boldsymbol{\nu}}+\lambda_{\boldsymbol{\nu}}(m)[\boldsymbol{\xi}, \boldsymbol{\nu}]+a \lambda_{\boldsymbol{\nu}}(m) D_{\boldsymbol{\nu}_{\perp}}\right] .
\end{aligned}
$$

Clearly, we can choose $a>0$ uniquely so that

$$
a \lambda_{\boldsymbol{\nu}}(m)\|\boldsymbol{\nu}\|^{2}=\rho_{m}\left(\left(i D_{\boldsymbol{\nu}_{\perp}}\right)_{M}(m),\left(i D_{\boldsymbol{\nu}_{\perp}}\right)_{M}(m)\right),
$$

so that $\left\langle\dot{\gamma}_{2}(0), \boldsymbol{\nu}_{\perp}\right\rangle=\left\langle\dot{\gamma}_{1}(0), \boldsymbol{\nu}_{\perp}\right\rangle$. Having fixed $a$, we can then choose $\dot{y}(0)$ uniquely so that

$$
\dot{y}(0)\|\boldsymbol{\nu}\|^{2}=\rho_{m}\left(\left(i D_{\boldsymbol{\nu}}\right)_{M}(m),\left(i D_{\boldsymbol{\nu}_{\perp}}\right)_{M}(m)\right),
$$

so that we also have $\left\langle\dot{\gamma}_{2}(0), \boldsymbol{\nu}\right\rangle=\left\langle\dot{\gamma}_{1}(0), \boldsymbol{\nu}\right\rangle$. Finally, if we set

$$
\boldsymbol{v}_{1}:=\left(\begin{array}{cc}
0 & 1 \\
-1 & 0
\end{array}\right), \quad \boldsymbol{v}_{2}:=\left(\begin{array}{cc}
0 & i \\
i & 0
\end{array}\right)
$$

we can choose $\boldsymbol{\xi} \in \operatorname{span}_{\mathbb{R}}\left\{\boldsymbol{v}_{1}, \boldsymbol{v}_{2}\right\}$ uniquely so that

$$
\lambda_{\boldsymbol{\nu}}(m)\left\langle[\boldsymbol{\xi}, i \boldsymbol{\nu}], \boldsymbol{v}_{j}\right\rangle=\rho_{m}\left(\left(\boldsymbol{v}_{j}\right)_{M}(m),\left(i D_{\boldsymbol{\nu}_{\perp}}\right)_{M}(m)\right),
$$

so that in addition $\left\langle\dot{\gamma}_{2}(0), \boldsymbol{v}_{j}\right\rangle=\left\langle\dot{\gamma}_{1}(0), \boldsymbol{v}_{j}\right\rangle$ for $j=1,2$. With this choices, $\gamma_{1}$ and $\gamma_{2}$ agree to first order at 0 .

Let us remark that when $\tau$ is sufficiently small $\gamma_{2}(\tau)$ has eigenvalues

$$
\lambda_{1}\left(\gamma_{2}(\tau)\right)=y(\tau)\left(\nu_{1}-a \tau \nu_{2}\right)>\lambda_{2}\left(\gamma_{2}(\tau)\right)=y(\tau)\left(\nu_{2}+a \tau \nu_{1}\right)
$$

Hence

$$
t(B, \boldsymbol{\nu}):=-\frac{a \tau}{\nu_{1}+\nu_{2}} \cdot \frac{\nu_{1}^{2}+\nu_{2}^{2}}{\nu_{1}-\nu_{2}+a \tau\left(\nu_{1}+\nu_{2}\right)} .
$$

Thus if $\nu_{1}+\nu_{2}>0$ then $\gamma_{2}(\tau) \notin \mathcal{B}_{\nu}$ when $\tau \in(0, \epsilon)$; since $\gamma_{1}$ and $\gamma_{2}$ agree to second order at 0 , we also have $\Phi_{G}(m+\tau \Upsilon(m)) \notin \mathcal{B}_{\nu}$ when $\tau \sim 0^{+}$.

The argument when $\nu_{1}+\nu_{2}<0$ is similar.
The proof of Theorem 1.3.2 is complete.

### 4.6 Proof of Theorem 1.3.3

### 4.6.1 An a priori polynomial bound

Let us record the following rough a priori polynomial bound.
Lemma 4.6.1. There is a constant $C_{\nu}>0$ such that for any $x \in X$ one has

$$
\left|\Pi_{k \nu}(x, x)\right| \leq C_{\nu} k^{d+1}
$$

for $k \gg 0$.
Proof of Lemma 4.6.1. Let $r: S^{1} \times X \rightarrow X$ be the standard structure action on the unit circle bundle $X$. Thus there is a natural decomposition

$$
H(X)=\bigoplus_{l=0}^{+\infty} H(X)_{l}
$$

into isotypes for the $S^{1}$-action.
Since $\tilde{\mu}$ commutes with the structure action $S^{1}$ on $X$, we have

$$
H(X)_{k \nu}=\bigoplus_{l=0}^{+\infty} H(X)_{k \nu} \cap H(X)_{l}
$$

On the other hand, by the theory of [GS82b] we have $H(X)_{k \nu} \cap H(X)_{l} \neq(0)$ only if the highest weight vector $\boldsymbol{r}(k \boldsymbol{\nu})$ of the representation indexed by $k \boldsymbol{\nu}$ satisfies

$$
\boldsymbol{r}(k \boldsymbol{\nu})=\left(k \nu_{1}-1, k \nu_{2}\right)=k \boldsymbol{\nu}+(-1,0) \in l \Phi_{G}(M) \subset \mathfrak{g} .
$$

Let us define

$$
a_{G}:=\min \left\|\Phi_{G}\right\|, \quad A_{G}:=\max \left\|\Phi_{G}\right\| .
$$

Thus $A_{G} \geq a_{G}>0$. Therefore, we need to have

$$
l a_{G} \leq\|\boldsymbol{r}(k \boldsymbol{\nu})\| \leq k\|\boldsymbol{\nu}\|+1 \Rightarrow l \leq L_{1}(k):=\left\lceil\frac{\|\boldsymbol{\nu}\|}{a_{G}} k+\frac{1}{a_{G}}\right\rceil .
$$

Similarly,

$$
k\|\boldsymbol{\nu}\|-1 \leq\|\boldsymbol{r}(k \boldsymbol{\nu})\| \leq l A_{G} \Rightarrow L_{2}(k):=\left\lfloor\frac{\|\boldsymbol{\nu}\|}{A_{G}} k-\frac{1}{A_{G}}\right\rfloor \leq l .
$$

On the other hand, in view of the asymptotic expansion of $\Pi_{k}(x, x)$ from the article [Zel98] we also have $\Pi_{l}(x, x) \leq 2(l / \pi)^{d}$ for $l \gg 0$. We conclude that

$$
\Pi_{k \nu}(x, x) \leq \sum_{l=L_{1}(k)}^{L_{2}(k)} \Pi_{l}(x, x) \leq \frac{2}{\pi^{d}} \sum_{l=L_{1}(k)}^{L_{2}(k)} l^{d} \leq C_{\nu} k^{d+1}
$$

for some $C_{\nu}>0$.

### 4.6.2 The proof

We shall use the following shorthand notation. If $x \in X, g \in G, t \in T$, let us set

$$
x(g, t):=\tilde{\mu}_{g t^{-1} g^{-1}}(x) ;
$$

similarly, if $m \in M$

$$
m(g, t):=\mu_{g t^{-1} g^{-1}}(m)
$$

If $t=e^{i \boldsymbol{\vartheta}}:=\left(e^{i \vartheta_{1}}, e^{i \vartheta_{2}}\right)$, we shall write $x(g, t)=x(g, \boldsymbol{\vartheta}), m(g, t)=m(g, \boldsymbol{\vartheta})$. Since $\tilde{\mu}$ is a lifting of $\mu$, if $m=\pi(x)$ then

$$
m(g, \boldsymbol{\vartheta})=\pi(x(g, \boldsymbol{\vartheta})) .
$$

Proof of Theorem 1.3.3. We replace $\boldsymbol{\nu}$ by $k \boldsymbol{\nu}$ in (3.9), and use the angular coordinates on $T$, we obtain

$$
\begin{equation*}
\Pi_{k \nu}(x, y)=\frac{k\left(\nu_{1}-\nu_{2}\right)}{(2 \pi)^{2}} \int_{-\pi}^{\pi} \int_{-\pi}^{\pi} e^{-i k\langle\boldsymbol{\nu}, \boldsymbol{\vartheta}\rangle} \Delta\left(e^{i \vartheta}\right) F\left(e^{i \vartheta} ; x, y\right) \mathrm{d} \boldsymbol{\vartheta} \tag{4.38}
\end{equation*}
$$

here $e^{i \vartheta}=\left(e^{i \vartheta_{1}}, e^{i \vartheta_{2}}\right)$.
For $\delta>0$, let us define

$$
V_{\delta}:=\left\{(x, y) \in X: \operatorname{dist}_{X}(x, G \cdot y) \geq \delta\right\}
$$

Proposition 4.6.1. For any $\delta>0$, we have $\Pi_{k \nu}(x, y)=O\left(k^{-\infty}\right)$ uniformly on $V_{\delta}$.

Proof of Proposition 4.6.1. The singular support of $\Pi$ is the diagonal in $X \times$ $X$. Therefore,

$$
\beta:((x, y), g T, t) \in V_{\delta} \times G / T \times T \mapsto \Pi(x(g, t), y) \in \mathbb{C}
$$

is $\mathcal{C}^{\infty}$. The same then holds of $((x, y), t) \in V_{\delta} \times T \mapsto \Delta(t) F(t ; x, y)$. Hence its Fourier transform (4.38) is rapidly decreasing for $k \rightarrow+\infty$.

We are thus reduced to assuming that $\operatorname{dist}_{X}(x, G \cdot y)<\delta$ for some fixed $\delta>0$. Let $\varrho \in \mathcal{C}_{0}^{\infty}(\mathbb{R})$ be $\equiv 1$ on $[-1,1]$ and $\equiv 0$ on $\mathbb{R} \backslash(-2,2)$. We can write

$$
\Pi_{\nu}(x, y)=\Pi_{\boldsymbol{\nu}}(x, y)_{1}+\Pi_{\boldsymbol{\nu}}(x, y)_{2}
$$

where the two summands on the right are defined as in (4.38), but with the integrand of $F(t ; x, y)$ multiplied respectively by $\varrho\left(\delta^{-1} \operatorname{dist}_{X}(x(g, \boldsymbol{\vartheta}), y)\right)$ and $1-\varrho\left(\delta^{-1} \operatorname{dist}_{X}(x(g, \boldsymbol{\vartheta}), y)\right)$.
Lemma 4.6.2. $\Pi_{k \nu}(x, y)_{2}=O\left(k^{-\infty}\right)$ for $k \rightarrow+\infty$.
Proof of Lemma 4.6.2. On the support of the integrand in $\Pi_{k \nu}(x, y)_{2}$ we have $\operatorname{dist}_{X}(x(g, t), y) \geq \delta$. We can then apply with minor changes the argument in the proof of Proposition 4.6.1.

On the support of the integrand in $\Pi_{k \nu}(x, y)_{1}, \operatorname{dist}_{X}(x(g, t), y) \leq 2 \delta ;$ therefore, perhaps after discarding a smoothing term contributing negligibly to the asymptotics, we can apply the explicit description of $\Pi$ as an FIO. With some passages, we obtain in place of (4.38):

$$
\begin{align*}
& \Pi_{k \nu}(x, y) \sim \Pi_{k \nu}(x, y)_{1}  \tag{4.39}\\
& \quad \sim \frac{k^{2}\left(\nu_{1}-\nu_{2}\right)}{(2 \pi)^{2}} \int_{-\pi}^{\pi} \int_{-\pi}^{\pi} \int_{G / T} \int_{0}^{+\infty} e^{i k \Psi_{x, y}} \mathcal{A}_{x, y} \mathrm{~d} u \mathrm{dV}_{G / T}(g T) \mathrm{d} \boldsymbol{\vartheta}
\end{align*}
$$

we have applied the rescaling $u \mapsto k u$ to the parameter in (4.38), and set

$$
\begin{gather*}
\Psi_{x, y}=\Psi_{x, y}(u, \boldsymbol{\vartheta}, g T):=u \psi\left(\tilde{\mu}_{g e^{-i \vartheta} g^{-1}}(x), y\right)-\langle\boldsymbol{\nu}, \boldsymbol{\vartheta}\rangle  \tag{4.40}\\
\mathcal{A}_{x, y}=\mathcal{A}_{x, y}(u, \boldsymbol{\vartheta}, g T):=\Delta\left(e^{i \vartheta}\right) s^{\prime}\left(\tilde{\mu}_{g e^{-i \vartheta} g^{-1}}(x), y, k u\right),
\end{gather*}
$$

with

$$
s^{\prime}\left(\tilde{\mu}_{g e^{-i \vartheta} \vartheta_{g^{-1}}}(x), y, k u\right):=s\left(\tilde{\mu}_{g e^{-i} \vartheta_{g^{-1}}}(x), y, k u\right) \varrho\left(\delta^{-1} \operatorname{dist}_{X}(x(g, \boldsymbol{\vartheta}), y)\right) .
$$

Lemma 4.6.3. Only a rapidly decreasing contribution to the asymptotics is lost, if in (4.39) integration in $\mathrm{d} u$ is restricted to an interval of the form $(1 / D, D)$ for some $D \gg 0$.

Proof of Lemma 4.6.3. Suppose that $x, y \in X,\left(g_{0} T, e^{i \vartheta_{0}}\right) \in(G / T) \times T$ and

$$
\operatorname{dist}_{X}\left(x\left(g_{0}, \boldsymbol{\vartheta}_{0}\right), y\right)<\delta
$$

In view of [SZ02], in any system of local coordinates we have

$$
\begin{equation*}
\mathrm{d}_{\left(x\left(g_{0}, \vartheta_{0}\right), y\right)} \psi=\left(\alpha_{x\left(g_{0}, \vartheta_{0}\right)},-\alpha_{y}\right)+O(\delta) . \tag{4.41}
\end{equation*}
$$

Let $\mathrm{d}^{(\boldsymbol{\vartheta})}$ denote the differential with respect to the variable $\boldsymbol{\vartheta}$. If $i \boldsymbol{\eta} \in \mathfrak{t}$, we obtain with $m_{x}:=\pi(x)$ :

$$
\begin{align*}
& \left.\frac{\mathrm{d}}{\mathrm{~d} \tau} x\left(g_{0}, \boldsymbol{\vartheta}_{0}+\tau \boldsymbol{\eta}\right)\right|_{\tau=0}=-\operatorname{Ad}_{g_{0}}(i \boldsymbol{\eta})_{X}\left(x\left(g_{0}, \boldsymbol{\vartheta}_{0}\right)\right)  \tag{4.42}\\
& \quad=-\operatorname{Ad}_{g_{0}}(i \boldsymbol{\eta})_{M}\left(x\left(g_{0}, \boldsymbol{\vartheta}_{0}\right)\right)^{\sharp}+\left\langle\Phi_{G}\left(m_{x}\left(g_{0}, \boldsymbol{\vartheta}_{0}\right)\right), \operatorname{Ad}_{g_{0}}(i \boldsymbol{\eta})\right\rangle \partial_{\theta} .
\end{align*}
$$

On the other hand, as $\Phi_{G}$ is $G$-equivariant we get

$$
\begin{align*}
& \left\langle\Phi_{G}\left(m_{x}\left(g_{0}, \boldsymbol{\vartheta}_{0}\right)\right), \operatorname{Ad}_{g_{0}}(i \boldsymbol{\eta})\right\rangle=\left\langle\operatorname{Ad}_{g_{0}^{-1}}\left(\Phi_{G}\left(m_{x}\left(g_{0}, \boldsymbol{\vartheta}_{0}\right)\right)\right), i \boldsymbol{\eta}\right\rangle  \tag{4.43}\\
& \quad=\left\langle\Phi_{G}\left(\tilde{\mu}_{g_{0}^{-1}}\left(m_{x}\left(g_{0}, \boldsymbol{\vartheta}_{0}\right)\right)\right), i \boldsymbol{\eta}\right\rangle=\left\langle\Phi_{T}\left(\tilde{\mu}_{g_{0}^{-1}}\left(m_{x}\left(g_{0}, \boldsymbol{\vartheta}_{0}\right)\right)\right), i \boldsymbol{\eta}\right\rangle .
\end{align*}
$$

Now, (4.41), (4.42) and (4.43) imply

$$
\begin{aligned}
\frac{\mathrm{d}}{\mathrm{~d} \tau} & \left.\psi\left(x\left(g_{0}, \boldsymbol{\vartheta}_{0}+\tau \boldsymbol{\eta}\right), y\right)\right|_{\tau=0} \\
& =-\mathrm{d}_{\left(x\left(g_{0}, \boldsymbol{\vartheta}_{0}\right), y\right)} \psi\left(\operatorname{Ad}_{g_{0}}(i \boldsymbol{\eta})_{X}\left(x\left(g_{0}, \boldsymbol{\vartheta}_{0}\right)\right), 0\right) \\
& =-\alpha_{x\left(g_{0}, \boldsymbol{\vartheta}_{0}\right)}\left(\operatorname{Ad}_{g_{0}}(i \boldsymbol{\eta})_{X}\left(x\left(g_{0}, \boldsymbol{\vartheta}_{0}\right)\right)\right)+\langle O(\delta), \boldsymbol{\eta}\rangle \\
& =\left\langle\operatorname{Ad}_{g_{0}^{-1}}\left(\frac{1}{i} \Phi_{T}\left(m_{x}\left(g_{0}, \boldsymbol{\vartheta}_{0}\right)\right)\right)+O(\delta), \boldsymbol{\eta}\right\rangle .
\end{aligned}
$$

Recalling (4.40), we obtain

$$
\begin{equation*}
\mathrm{d}_{\left(u, g_{0} T, \vartheta_{0}\right)}^{(\vartheta)} \Psi_{x, y}=\frac{u}{i} \Phi_{T}\left(\mu_{g_{0}^{-1}}\left(m_{x}\right)\right)-\boldsymbol{\nu}+O(\delta) . \tag{4.44}
\end{equation*}
$$

By assumption, $\mathbf{0} \notin \Phi_{T}(M)$. Let us set

$$
a_{T}:=\min \left\|\Phi_{T}\right\|, \quad A_{T}:=\max \left\|\Phi_{T}\right\|
$$

Then $A_{T} \geq a_{T}>0$, and (4.44) implies

$$
\left\|\mathrm{d}_{\left(u, g_{0} T, \vartheta_{0}\right)}^{(\boldsymbol{\vartheta})} \Psi_{x, y}\right\| \geq \max \left\{u a_{T}-\|\boldsymbol{\nu}\|+O(\delta),\|\boldsymbol{\nu}\|-u A_{T}+O(\delta)\right\} .
$$

Thus if $D \gg 0$ and $u \geq D$ we have

$$
\begin{equation*}
\left\|\mathrm{d}_{\left(u, g_{0} T, \vartheta_{0}\right)}^{(\boldsymbol{\vartheta})} \Psi_{x, y}\right\| \geq \frac{a_{T}}{2} u+1, \tag{4.45}
\end{equation*}
$$

while for $0<u<1 / D$

$$
\begin{equation*}
\left\|\mathrm{d}_{\left(u, g_{0} T, \boldsymbol{\vartheta}_{0}\right)}^{(\boldsymbol{\vartheta})} \Psi_{x, y}\right\| \geq \frac{\|\boldsymbol{\nu}\|}{2} . \tag{4.46}
\end{equation*}
$$

The Lemma then follows from (4.45) and (4.46) by a standard iterated integration by parts in $\boldsymbol{\vartheta}$ (in view of the compactness of $T$ ).

Suppose that $\varrho \in \mathcal{C}_{0}^{\infty}((0,+\infty))$ is $\equiv 1$ on $(1 / D, D)$ and is supported on $(1 /(2 D), 2 D)$. By Lemma 4.6.3, the asymptotics of (4.39) is unaltered, if the integrand is multiplied by $\varrho(u)$. Thus we obtain

$$
\begin{equation*}
\Pi_{k \nu}(x, y) \sim \frac{k^{2}\left(\nu_{1}-\nu_{2}\right)}{(2 \pi)^{2}} \int_{(-\pi, \pi)^{2}} \int_{G / T} \int_{1 / 2 D}^{2 D} e^{i k \Psi_{x, y}} \mathcal{A}_{x, y}^{\prime} \mathrm{d} u \mathrm{~d} V_{G / T}(g T) \mathrm{d} \boldsymbol{\vartheta} \tag{4.47}
\end{equation*}
$$

where we have set

$$
\mathcal{A}_{x, y}^{\prime}(u, \boldsymbol{\vartheta}, g T):=\varrho(u) \mathcal{A}_{x, y}(u, \boldsymbol{\vartheta}, g T)
$$

and the integration in $\mathrm{d} u$ is now over a compact interval.
Let $\Im(z)$ denote the imaginary part of $z \in \mathbb{C}$. In view of Corollary 1.3 of [BdMS76], there exists a fixed constant $D^{\prime}$, depending only on $X$, such that

$$
\begin{equation*}
\Im\left(\psi\left(x^{\prime}, x^{\prime \prime}\right)\right) \geq D^{\prime} \operatorname{dist}_{X}\left(x^{\prime}, x^{\prime \prime}\right)^{2} \quad\left(x^{\prime}, x^{\prime \prime} \in X\right) \tag{4.48}
\end{equation*}
$$

Proposition 4.6.2. Uniformly for

$$
\begin{equation*}
\operatorname{dist}_{X}(x, G \cdot y) \geq C k^{\epsilon-1 / 2} \tag{4.49}
\end{equation*}
$$

we have $\Pi_{k \nu}(x, y)=O\left(k^{-\infty}\right)$.
Proof of Proposition 4.6.2. In the range (4.49), we have

$$
\operatorname{dist}_{X}(x(g, \boldsymbol{\vartheta}), y) \geq C k^{\epsilon-1 / 2}
$$

for every $g T \in G / T$ and $e^{i \vartheta} \in T$. In view of (4.40) and (4.48),

$$
\begin{align*}
\left|\partial_{u} \Psi_{x, y}(u, \boldsymbol{\vartheta}, g T)\right| & =|\psi(x(g, \boldsymbol{\vartheta}), y)| \geq \Im(\psi(x(g, \boldsymbol{\vartheta}), y))  \tag{4.50}\\
& \geq D \operatorname{dist}_{X}(x(g, \boldsymbol{\vartheta}), y)^{2} \geq D C^{2} k^{2 \epsilon-1} \tag{4.51}
\end{align*}
$$

Let us use the identity

$$
\begin{equation*}
-\frac{i}{k} \psi(x(g, \boldsymbol{\vartheta}), y)^{-1} \frac{\mathrm{~d}}{\mathrm{~d} u} e^{i k \Psi_{x, y}}=e^{i k \Psi_{x, y}} \tag{4.52}
\end{equation*}
$$

to iteratively integrate by parts in $\mathrm{d} u$ in (4.47); then by (4.50) at each step we introduce a factor $O\left(k^{-2 \epsilon}\right)$. The claim follows.

To complete the proof of Theorem 1.3.3, we need to establish the following.
Proposition 4.6.3. Uniformly for

$$
\begin{equation*}
\operatorname{dist}_{X}\left(x, G \cdot X_{\nu}^{T}\right) \geq C k^{\epsilon-1 / 2} \tag{4.53}
\end{equation*}
$$

we have $\Pi_{k \nu}(x, y)=O\left(k^{-\infty}\right)$.
Remark 4. Let $\operatorname{dist}_{M}$ denote the distance function on $M$; if $m=\pi(x)$, then $\operatorname{dist}_{X}\left(x, G \cdot X_{\nu}^{T}\right)=\operatorname{dist}_{M}\left(m, G \cdot M_{\nu}^{T}\right)$.

Proof of Proposition 4.6.3. Since $G$ acts on $M$ as a group of Riemannian isometries, (4.53) means that for any $g \in G$ we have

$$
C k^{\epsilon-1 / 2} \leq \operatorname{dist}_{M}\left(m, \mu_{g}\left(M_{\nu}^{T}\right)\right)=\operatorname{dist}_{M}\left(\mu_{g^{-1}}(m), M_{\nu}^{T}\right)
$$

On the other hand, as $-i \Phi_{T}$ is transverse to $\mathbb{R}_{+} \cdot \boldsymbol{\nu}$, by the discussion in §2.1.3 of [Pao17] there is a constant $b_{\boldsymbol{\nu}}>0$ such that every $u \in[1 /(2 D), 2 D]$ we have

$$
\begin{equation*}
\left\|-i u \Phi_{T}\left(\mu_{g^{-1}}(m)\right)-\boldsymbol{\nu}\right\| \geq b_{\nu} C k^{\epsilon-1 / 2} \tag{4.54}
\end{equation*}
$$

Let us consider (4.47) with $x=y$ :

$$
\begin{align*}
& \Pi_{k \nu}(x, x) \sim  \tag{4.55}\\
& \quad \frac{k^{2}\left(\nu_{1}-\nu_{2}\right)}{(2 \pi)^{2}} \int_{-\pi}^{\pi} \int_{-\pi}^{\pi} \int_{G / T} \int_{1 / 2 D}^{2 D} e^{i k \Psi_{x, x}} \mathcal{A}_{x, x}^{\prime} \mathrm{d} u \mathrm{~d} V_{G / T}(g T) \mathrm{d} \boldsymbol{\vartheta}
\end{align*}
$$

Let us choose $\epsilon^{\prime} \in(0, \epsilon)$ and multiply the integrand in (4.55) by the identity

$$
\varrho\left(k^{1 / 2-\epsilon^{\prime}} \operatorname{dist}_{X}(x(g, \boldsymbol{\vartheta}), x)\right)+\left[1-\varrho\left(k^{1 / 2-\epsilon^{\prime}} \operatorname{dist}_{X}(x(g, \boldsymbol{\vartheta}), x)\right)\right]=1 .
$$

Here $\varrho$ is as in the discussion preceding Lemma 4.6.2. We obtain a further splitting

$$
\begin{equation*}
\Pi_{k \nu}(x, x) \sim \Pi_{k \nu}(x, x)_{a}+\Pi_{k \nu}(x, x)_{b} \tag{4.56}
\end{equation*}
$$

where $\Pi_{k \nu}(x, x)_{a}$ is given by (4.55) with $\mathcal{A}_{x, x}^{\prime}$ replaced by

$$
\begin{equation*}
\mathcal{B}_{x, x}^{\prime}:=\varrho\left(k^{1 / 2-\epsilon^{\prime}} \operatorname{dist}_{X}(x(g, \boldsymbol{\vartheta}), x)\right) \mathcal{A}_{x, x}^{\prime} \tag{4.57}
\end{equation*}
$$

similarly $\Pi_{k \nu}(x, x)_{b}$ is given by (4.55) with $\mathcal{A}_{x, x}^{\prime}$ replaced by

$$
\mathcal{B}_{x, x}^{\prime \prime}:=\left[1-\varrho\left(k^{1 / 2-\epsilon^{\prime}} \operatorname{dist}_{X}(x(g, \boldsymbol{\vartheta}), x)\right)\right] \mathcal{A}_{x, x}^{\prime}
$$

Lemma 4.6.4. $\Pi_{k \nu}(x, x)_{b}=O\left(k^{-\infty}\right)$ as $k \rightarrow+\infty$.
Proof of Lemma 4.6.4. On the support of $\mathcal{B}_{x, x}^{\prime \prime}$, we have

$$
\operatorname{dist}_{X}(x(g, \boldsymbol{\vartheta}), x) \geq k^{\epsilon^{\prime}-1 / 2}
$$

Thus we may again appeal to (4.52) and iteratively integrate by parts in $\mathrm{d} u$, introducing at each step a factor $O\left(k^{-1} k^{1-2 \epsilon^{\prime}}\right)=O\left(k^{-2 \epsilon^{\prime}}\right)$.

Thus the proof of the Theorem will be complete once we establish the following.

Lemma 4.6.5. $\Pi_{k \nu}(x, x)_{a}=O\left(k^{-\infty}\right)$ as $k \rightarrow+\infty$.
Before attacking the proof of Lemma 4.6.5, let us prove the following.
Lemma 4.6.6. If (4.53) holds, then for any $u \in[1 /(2 D), 2 D]$ and $k \gg 0$

$$
\left\|\mathrm{d}_{(u, g T, \vartheta)}^{\vartheta} \Psi_{x, x}\right\| \geq \frac{b_{\nu}}{2} C k^{\epsilon-1 / 2}
$$

on the support of $\mathcal{B}_{x, x}^{\prime}$.

Proof of Lemma 4.6.6. On the support of $\mathcal{B}_{x, x}^{\prime}$, we have

$$
\operatorname{dist}_{X}(x(g, \boldsymbol{\vartheta}), x) \leq 2 k^{\epsilon^{\prime}-1 / 2}
$$

Thus, instead of (4.41) we have

$$
\mathrm{d}_{(x(g, \vartheta), x)} \psi=\left(\alpha_{x(g, \vartheta)},-\alpha_{x}\right)+O\left(k^{\epsilon^{\prime}-1 / 2}\right)
$$

Therefore, in place of (4.44), on the support of $\mathcal{B}_{x, x}^{\prime}$ we have

$$
\mathrm{d}_{(u, g T, \boldsymbol{\vartheta})}^{(\boldsymbol{\vartheta})} \Psi_{x, x}=\frac{u}{i} \Phi_{T}\left(\mu_{g^{-1}}\left(m_{x}\right)\right)-\boldsymbol{\nu}++O\left(k^{\epsilon^{\prime}-1 / 2}\right)
$$

Thus in view of (4.54) the claim follows from $0<\epsilon^{\prime}<\epsilon$.
Given Lemma 4.6.6, we can prove Lemma 4.6.5 essentially by iteratively integrating by parts in $\mathrm{d} \boldsymbol{\vartheta}$.

Proof of Lemma 4.6.5. Since $\tilde{\mu}$ is free on $X_{\mathcal{O}_{\nu}}^{G}$, it is also free on a small tubular neighbourhood $X^{\prime}$ of $X_{\mathcal{O}_{\nu}}^{G}$ in $X$. Without loss, we may restrict our analysis to $X$ in view of Theorem 1.3.3. On the support of $\mathcal{B}_{x, x}^{\prime}$, therefore, $e^{i \vartheta} \in T$ varies in a small neighbourhood of $I_{2}$. Let $f: T \rightarrow[0,+\infty)$ be a bump function compactly supported in a small neighbourhood $U \subset T$ of $I_{2}$ (identified with $(1,1)$ ), and identically $\equiv 1$ near $I_{2}$. Then we obtain

$$
\begin{aligned}
& \Pi_{k \nu}(x, x)_{a} \sim \\
& \quad \frac{k^{2}\left(\nu_{1}-\nu_{2}\right)}{(2 \pi)^{2}} \int_{U} \int_{G / T} \int_{1 / 2 D}^{2 D} e^{i k \Psi_{x, x}} f(t) \mathcal{B}_{x, x}^{\prime} \mathrm{d} u \mathrm{dV}_{G / T}(g T) \mathrm{d} \boldsymbol{\vartheta},
\end{aligned}
$$

Let us introduce the differential operator

$$
P=\sum_{h=1}^{2} \frac{\partial_{\vartheta_{h}} \Psi_{x, x}}{\left(\partial_{\vartheta_{1}} \Psi_{x, x}\right)^{2}+\left(\partial_{\vartheta_{2}} \Psi_{x, x}\right)^{2}} \frac{\partial}{\partial \vartheta_{h}},
$$

so that

$$
\frac{1}{i k} P\left(e^{i k \Psi_{x, x}}\right)=e^{i k \Psi_{x, x}}
$$

Thus,

$$
\begin{aligned}
& \int_{U} e^{i k \Psi_{x, x}} f(t) \mathcal{B}_{x, x}^{\prime} \mathrm{d} \boldsymbol{\vartheta} \\
& \quad=\frac{1}{i k} \sum_{h=1}^{2} \int_{U} \frac{\partial_{\vartheta_{h}} \Psi_{x, x}}{\left(\partial_{\vartheta_{1}} \Psi_{x, x}\right)^{2}+\left(\partial_{\vartheta_{2}} \Psi_{x, x}\right)^{2}} \frac{\partial}{\partial \vartheta_{h}}\left[e^{i k \Psi_{x, x}}\right] f\left(e^{i \boldsymbol{\vartheta}}\right) \mathcal{B}_{x, x}^{\prime} \mathrm{d} \boldsymbol{\vartheta} \\
& \quad=\frac{i}{k} \sum_{h=1}^{2} \int_{U} e^{i k \Psi_{x, x}} \frac{\partial}{\partial \vartheta_{h}}\left[\frac{\partial_{\vartheta_{h}} \Psi_{x, x}}{\left(\partial_{\vartheta_{1}} \Psi_{x, x}\right)^{2}+\left(\partial_{\vartheta_{2}} \Psi_{x, x}\right)^{2}} f\left(e^{i \boldsymbol{\vartheta}}\right) \mathcal{B}_{x, x}^{\prime}\right] \mathrm{d} \boldsymbol{\vartheta} \\
& \\
& =\frac{i}{k} \int_{U} e^{i k \Psi_{x, x}} P^{t}\left[f\left(e^{i \boldsymbol{\vartheta}}\right) \mathcal{B}_{x, x}^{\prime}\right] \mathrm{d} \boldsymbol{\vartheta}
\end{aligned}
$$

where

$$
\begin{equation*}
P^{t}(\gamma):=\sum_{h=1}^{2} \frac{\partial}{\partial \vartheta_{h}}\left[\frac{\partial_{\vartheta_{h}} \Psi_{x, x}}{\left(\partial_{\vartheta_{1}} \Psi_{x, x}\right)^{2}+\left(\partial_{\vartheta_{2}} \Psi_{x, x}\right)^{2}} \gamma\right] \tag{4.58}
\end{equation*}
$$

Iterating, for any $r \in \mathbb{N}$ we have

$$
\begin{equation*}
\int_{U} e^{i k \Psi_{x, x}} f\left(e^{i \boldsymbol{\vartheta}}\right) \mathcal{B}_{x, x}^{\prime} \mathrm{d} \boldsymbol{\vartheta}=\frac{i^{r}}{k^{r}} \int_{U} e^{i k \Psi_{x, x}}\left(P^{t}\right)^{r}\left[f\left(e^{i \boldsymbol{\vartheta}}\right) \mathcal{B}_{x, x}^{\prime}\right] \mathrm{d} \boldsymbol{\vartheta} \tag{4.59}
\end{equation*}
$$

Let us consider the function

$$
\mathcal{D}: \boldsymbol{\vartheta} \mapsto \operatorname{dist}_{X}(x(g, \boldsymbol{\vartheta}))=\operatorname{dist}_{X}\left(\tilde{\mu}_{e^{-i} \vartheta} \circ \tilde{\mu}_{g^{-1}}(x), \tilde{\mu}_{g^{-1}}(x)\right) .
$$

We have the following.
Lemma 4.6.7. For $\boldsymbol{\vartheta} \sim \mathbf{0}$, we have

$$
\operatorname{dist}_{X}(x(g, \boldsymbol{\vartheta}), x)=F_{1}(g T, \boldsymbol{\vartheta})+F_{2}(g T, \boldsymbol{\vartheta})+\ldots,
$$

where $F_{j}(g T, \boldsymbol{\vartheta})$ is homogeneous of degree $j$ in $\boldsymbol{\vartheta}$, and $\mathcal{C}^{\infty}$ for $\boldsymbol{\vartheta} \neq \mathbf{0}$. In addition, $F_{1}(g T, \boldsymbol{\vartheta})=\left\|\operatorname{Ad}_{g}(\boldsymbol{\vartheta})_{X}(x)\right\|=\left\|\boldsymbol{\vartheta}_{X}\left(\tilde{\mu}_{g^{-1}}(x)\right)\right\|$.

For any $c \in \mathbb{N}$ let $\mathcal{D}^{(c)}$ denote a generic iterated derivative of the form

$$
\frac{\partial^{c} \mathcal{D}}{\partial \vartheta_{i_{1}} \cdots \partial \vartheta_{i_{c}}} ;
$$

clearly $\mathcal{D}^{(c)}$ is not uniquely determined by $c$. By Lemma 4.6.7, as $k \rightarrow+\infty$

$$
\mathcal{D}^{(c)}\left(\frac{\boldsymbol{\vartheta}}{\sqrt{k}}\right)=O\left(k^{(c-1)\left(1 / 2-\epsilon^{\prime}\right)}\right)
$$

where $\varrho\left(k^{1 / 2-\epsilon^{\prime}} \mathcal{D}\right) \neq 1$. For any multi-index $\boldsymbol{C}=\left(c_{1}, \cdots, c_{s}\right)$ let us denote by $\mathcal{D}^{(C)}$ a generic product of the form $\mathcal{D}^{\left(c_{1}\right)} \cdots \mathcal{D}^{\left(c_{s}\right)}$; then

$$
\mathcal{D}^{(C)}\left(\frac{\boldsymbol{\vartheta}}{\sqrt{k}}\right)=O\left(k^{\left(1 / 2-\epsilon^{\prime}\right) \sum_{j}\left(c_{j}-1\right)}\right) .
$$

Lemma 4.6.8. For any $r \in \mathbb{N},\left(P^{t}\right)^{r}\left(f(t) \mathcal{B}_{x, x}^{\prime}\right)$ is a linear combination of summands of the form

$$
\begin{equation*}
\varrho^{(b)}\left(k^{1 / 2-\epsilon^{\prime}} \mathcal{D}_{k}(\boldsymbol{\vartheta})\right) \frac{P_{a_{1}}\left(\Psi_{x, x}, \partial \Psi_{x, x}\right)}{\left[\left(\partial_{\vartheta_{1}} \Psi_{x, x}\right)^{2}+\left(\partial_{\vartheta_{2}} \Psi_{x, x}\right)^{2}\right]^{a_{2}}} k^{b\left(1 / 2-\epsilon^{\prime}\right)} \mathcal{D}^{(C)}, \tag{4.60}
\end{equation*}
$$

times omitted factors bounded in $k$ depending on $f$ and its derivatives, where

1. $P_{a_{1}}$ denotes a generic differential polynomial in $\Psi_{x, x}$ homogeneous of degree $a_{1}$ in the first derivatives $\partial \Psi_{x, x}$;
2. If $a:=2 a_{2}-a_{1}$, then $a, b, C$ are subjected to the bound

$$
\begin{equation*}
a+b+\sum_{j=1}^{r}\left(c_{j}-1\right) \leq 2 r \tag{4.61}
\end{equation*}
$$

(the sum is over $c_{j}>0$ );
3. $\boldsymbol{C}$ is not zero if and only if $b>0$.

Here $\varrho^{(b)}$ is the $b$-th derivative of the one-variable real function $\varrho$.
Proof of Lemma 4.6.8. Let us set $F:=f\left(e^{i \vartheta}\right) \mathcal{B}_{x, x}^{\prime}$. For $r=1$, we have

$$
\begin{align*}
& \frac{\partial}{\partial \vartheta_{h}}\left[\frac{\partial_{\vartheta_{h}} \Psi_{x, x}}{\left(\partial_{\vartheta_{1}} \Psi_{x, x}\right)^{2}+\left(\partial_{\vartheta_{2}} \Psi_{x, x}\right)^{2}} F\right]  \tag{4.62}\\
& \quad=\frac{\partial_{\vartheta_{h}} \Psi_{x, x}}{\left(\partial_{\vartheta_{1}} \Psi_{x, x}\right)^{2}+\left(\partial_{\vartheta_{2}} \Psi_{x, x}\right)^{2}} \frac{\partial F}{\partial \vartheta_{h}}+F \frac{\partial}{\partial \vartheta_{h}}\left[\frac{\partial_{\vartheta_{h}} \Psi_{x, x}}{\left(\partial_{\vartheta_{1}} \Psi_{x, x}\right)^{2}+\left(\partial_{\vartheta_{2}} \Psi_{x, x}\right)^{2}}\right]
\end{align*}
$$

We have

$$
\begin{aligned}
& \frac{\partial_{\theta_{h}} \Psi_{x, x}}{\left(\partial_{\theta_{1}} \Psi_{x, x}\right)^{2}+\left(\partial_{\theta_{2}} \Psi_{x, x}\right)^{2}} \frac{\partial F}{\partial \vartheta_{h}} \\
& \quad=\frac{\partial_{\theta_{h}} \Psi_{x, x}}{\left(\partial_{\theta_{1}} \Psi_{x, x}\right)^{2}+\left(\partial_{\theta_{2}} \Psi_{x, x}\right)^{2}}\left[\frac{\partial f}{\partial \vartheta_{h}} \mathcal{B}_{x, x}^{\prime}+\frac{\partial \mathcal{B}_{x, x}^{\prime}}{\partial \vartheta_{h}} f\right] .
\end{aligned}
$$

Thus, in view of (4.57), the first summand on the right hand side of (4.62) splits as a linear combination of terms as in the statement, with $a_{1}=a_{2}=1, b$ and $\boldsymbol{C}$ both zero, or $a_{1}=a_{2}=1, b=1, \boldsymbol{C}=(1)$. Hence $a+b+\sum_{j}\left(c_{j}-1\right)=2$ in either case. On the other hand, the second summand on the right hand side of (4.62) satisfies

$$
\begin{aligned}
& F \frac{\partial}{\partial \vartheta_{h}}\left[\frac{\partial_{\vartheta_{h}} \Psi_{x, x}}{\left(\partial_{\vartheta_{1}} \Psi_{x, x}\right)^{2}+\left(\partial_{\vartheta_{2}} \Psi_{x, x}\right)^{2}}\right]=\frac{F}{\left[\left(\partial_{\vartheta_{1}} \Psi_{x, x}\right)^{2}+\left(\partial_{\vartheta_{2}} \Psi_{x, x}\right)^{2}\right]^{2}} \\
& \cdot\left\{\partial_{\vartheta_{h}, \vartheta_{h}}^{2} \Psi_{x, x}\left[\left(\partial_{\vartheta_{1}} \Psi_{x, x}\right)^{2}+\left(\partial_{\vartheta_{2}} \Psi_{x, x}\right)^{2}\right]-2 \partial_{\vartheta_{h}} \Psi_{x, x} \sum_{a=1}^{2} \partial_{\vartheta_{a}} \Psi_{x, x} \partial_{\vartheta_{a} \vartheta_{h}}^{2} \Psi_{x, x}\right\}
\end{aligned}
$$

This is of the stated type with $a_{1}=a_{2}=2, b$ and $\boldsymbol{C}$ both zero. Hence $a=4-2=2$.

Passing to the inductive step, let us consider (4.58) with $\gamma$ given by (4.60), and assume that (4.61) is satisfied. Let us write $\varrho^{(l)}$ for the factor in front in (4.60). We obtain a linear combination of expressions of the form

$$
\begin{equation*}
\frac{\partial}{\partial \vartheta_{h}}\left[\varrho^{(b)} \frac{P_{a_{1}+1}\left(\Psi_{x, x}, \partial \Psi_{x, x}\right)}{\left[\left(\partial_{\vartheta_{1}} \Psi_{x, x}\right)^{2}+\left(\partial_{\vartheta_{2}} \Psi_{x, x}\right)^{2}\right]^{a_{2}+1}} k^{b\left(1 / 2-\epsilon^{\prime}\right)} \mathcal{D}^{(C)} k^{b\left(1 / 2-\epsilon^{\prime}\right)} \mathcal{D}^{(C)}\right] \tag{4.63}
\end{equation*}
$$

It is now clear that (4.63) splits as a linear combination of summands of the following forms:

$$
\varrho^{(b)} \frac{P_{a^{\prime}}\left(\Psi_{x, x}, \partial \Psi_{x, x}\right)}{\left[\left(\partial_{\vartheta_{1}} \Psi_{x, x}\right)^{2}+\left(\partial_{\vartheta_{2}} \Psi_{x, x}\right)^{2}\right]^{a_{2}+1}} k^{b\left(1 / 2-\epsilon^{\prime}\right)} \mathcal{D}^{(C)} k^{b\left(1 / 2-\epsilon^{\prime}\right)} \mathcal{D}^{(C)},
$$

with $a^{\prime} \in\left\{a_{1}, a_{1}+1, a_{1}+2\right\}$;

$$
\begin{gathered}
\varrho^{(b)} \frac{P_{a_{1}+2}\left(\Psi_{x, x}, \partial \Psi_{x, x}\right)}{\left[\left(\partial_{\vartheta_{1}} \Psi_{x, x}\right)^{2}+\left(\partial_{\vartheta_{2}} \Psi_{x, x}\right)^{2}\right]^{a_{2}+2}} k^{b\left(1 / 2-\epsilon^{\prime}\right)} \mathcal{D}^{(C)} k^{b\left(1 / 2-\epsilon^{\prime}\right)} \mathcal{D}^{(C)} \\
\varrho^{(b+1)} \frac{P_{a_{1}+1}\left(\Psi_{x, x}, \partial \Psi_{x, x}\right)}{\left[\left(\partial_{\vartheta_{1}} \Psi_{x, x}\right)^{2}+\left(\partial_{\vartheta_{2}} \Psi_{x, x}\right)^{2}\right]^{a_{2}+1}} k^{(b+1)\left(1 / 2-\epsilon^{\prime}\right)} \mathcal{D}^{(C)} k^{b\left(1 / 2-\epsilon^{\prime}\right)} \mathcal{D}^{\left(C^{\prime}\right)},
\end{gathered}
$$

where $\boldsymbol{C}^{\prime}$ is of the form $\boldsymbol{C}^{\prime}=(1, \boldsymbol{C})$;

$$
\varrho^{(b)} \frac{P_{a_{1}+1}\left(\Psi_{x, x}, \partial \Psi_{x, x}\right)}{\left[\left(\partial_{\vartheta_{1}} \Psi_{x, x}\right)^{2}+\left(\partial_{\vartheta_{2}} \Psi_{x, x}\right)^{2}\right]^{a_{2}+1}} k^{b\left(1 / 2-\epsilon^{\prime}\right)} \mathcal{D}^{\left(C^{\prime}\right)}
$$

where $\boldsymbol{C}^{\prime}$ is obtained from $\boldsymbol{C}$ (if the latter is not zero) by replacing one of the $c_{j}$ 's by $c_{j}+1$, and leaving all the others unchanged.

In all these cases we obtain a term of the form (4.60), satisfying (4.61) with $r$ replaced by $r+1$. This completes the proof of Lemma 4.6.8.

As $0<\epsilon^{\prime}<\epsilon$, the general summand (4.61) is

$$
\begin{aligned}
O\left(k^{a(1 / 2-\epsilon)+\left[b+\sum_{j}\left(c_{j}-1\right)\right]\left(1 / 2-1 \epsilon^{\prime}\right)}\right) & =O\left(k^{\left[a+b+\sum_{j}\left(c_{j}-1\right)\right]\left(1 / 2-1 \epsilon^{\prime}\right)}\right) \\
& =O\left(k^{2 r\left(1 / 2-1 \epsilon^{\prime}\right)}\right)=O\left(k^{r\left(1-2 \epsilon^{\prime}\right)}\right) .
\end{aligned}
$$

Making use of the latter estimate in (4.59), we obtain the following
Corollary 4.6.1. For any $r \in \mathbb{N}$,

$$
\int_{U_{j}} e^{i k \Psi_{x, x}} f(t) \mathcal{B}_{x, x}^{\prime} \mathrm{d} \boldsymbol{\vartheta}=O\left(k^{-2 r \epsilon^{\prime}}\right)
$$

The proof of Lemma 4.6.5 is thus complete.
Given (4.56), Proposition 4.6.3 follows from Lemmata 4.6.4 and 4.6.5.
Thus the statement of Theorem 1.3.3 holds true when $x=y$. The general case follows from this and the Schwartz inequality

$$
\left|\Pi_{\nu}(x, y)\right| \leq \sqrt{\Pi_{\nu}(x, x)} \sqrt{\Pi_{\nu}(y, y)}
$$

in fact both factors on the right hand side have at most polynomial growth in $k$ by Lemma 4.6.1, and if say (4.53) holds, then the first one is rapidly decreasing. The proof of Theorem 1.3.3 is complete.

### 4.7 Proof of Theorem 1.3.4, 1.3.6 and 1.3.7

### 4.7.1 Preliminaries on local rescaled asymptotics

In the proof of Theorem 1.3.4, 1.3.6 and 1.3.7, we are interested in the asymptotics of $\Pi_{k \nu}\left(x^{\prime}, x^{\prime \prime}\right)$ when $\left(x^{\prime}, x^{\prime \prime}\right)$ approaches the diagonal of $X_{\mathcal{O}_{\nu}}^{G}$ in $X \times X$ along appropriate directions and at a suitable pace.

In Theorems 1.3.4 and 1.3.7, we consider $x^{\prime}=x^{\prime \prime}$ in a shrinking "onesided" neighbourhood of $X_{\mathcal{O}_{\nu}}^{G}$. In Theorem 1.3.6, we shall assume that
$\left(x^{\prime}, x^{\prime \prime}\right)$ approaches the diagonal in $X_{\mathcal{O}_{\nu}}^{G}$ along "horizontal" directions orthogonal to the orbits. We shall treat the former case in detail, and then briefly discuss the necessary changes for the latter. Suppose $x \in X_{\mathcal{O}_{\nu}}^{G}$ and let $m=\pi(x)$. Let us choose a system of HLC centered at $x$, and let us consider the collection of points

$$
x_{\tau, k}:=x+\frac{\tau}{\sqrt{k}} \Upsilon_{\nu}(m)
$$

where $k=1,2, \ldots$, and $|\tau| \leq C k^{\epsilon}$ for some fixed $C>0$ and $\epsilon \in(0,1 / 6)$. The sign of $\tau$ is chosen so that $\tau \Upsilon_{\nu}(m)$ is either zero or outer oriented. Thus $\tau\left(\nu_{1}+\nu_{2}\right) \geq 0$. We shall provide an integral expression for the asymptotics of $\Pi_{k \nu}\left(x_{\tau, k}, x_{\tau, k}\right)$ when $k \rightarrow+\infty$. Applying as before the Weyl integration and character formulae, inserting the micro-local description of $\Pi$ as an FIO and making use of the rescaling $u \mapsto k u, \boldsymbol{\vartheta} \mapsto \boldsymbol{\vartheta} / \sqrt{k}$, we obtain that, as $k \rightarrow+\infty$,

$$
\begin{align*}
& \Pi_{k \nu}\left(x_{\tau, k}, x_{\tau, k}\right) \sim  \tag{4.64}\\
& \quad \frac{k\left(\nu_{1}-\nu_{2}\right)}{(2 \pi)^{2}} \int_{G / T} \mathrm{~d} V_{G / T}(g T) \int_{-\infty}^{+\infty} \mathrm{d} \vartheta_{1} \int_{-\infty}^{+\infty} \mathrm{d} \vartheta_{2} \int_{0}^{+\infty} \mathrm{d} u \\
& {\left[e^{i k\left[u \psi\left(\tilde{\mu}_{g e^{-i} \vartheta^{-} / \sqrt{k}_{g^{-1}}}\left(x_{\tau, k}\right), x_{\tau, k}\right)-\langle\vartheta, \nu\rangle / \sqrt{k}\right]}\right.} \\
& \left.\Delta\left(e^{i \vartheta / \sqrt{k}}\right) s\left(\tilde{\mu}_{g e^{-i \vartheta / \sqrt{k}} g^{-1}}\left(x_{\tau, k}\right), x_{\tau, k}, k u\right)\right] .
\end{align*}
$$

Integration in $\boldsymbol{\vartheta}=\left(\vartheta_{1}, \vartheta_{2}\right)$ is over a ball centered at the origin and radius $O\left(k^{\epsilon}\right)$ in $\mathbb{R}^{2}$. A cut-off function of the form $\varrho\left(k^{\epsilon} \boldsymbol{\vartheta}\right)$ is implicitly incorporated into the amplitude.

In order to express the previous phase more explicitly, we need the following.
Definition 4.7.1. Let us define $\boldsymbol{\rho}=\boldsymbol{\rho}_{m}: G / T \rightarrow \mathfrak{t} \cong \mathbb{R}^{2}, g T \mapsto \boldsymbol{\rho}_{g T}$, by requiring

$$
\left\langle\boldsymbol{\rho}_{g T}, \boldsymbol{\vartheta}\right\rangle=\omega_{m}\left(\operatorname{Ad}_{g}\left(i D_{\vartheta}\right)_{M}(m), \Upsilon_{\boldsymbol{\nu}}(m)\right) \quad\left(\boldsymbol{\vartheta} \in \mathbb{R}^{2}\right) .
$$

Next, let the symmetric and positive definite matrix $E(g T)=E_{x}(g T)$ be defined by the equality

$$
\boldsymbol{\vartheta}^{t} E(g T) \boldsymbol{\vartheta}=\left\|\operatorname{Ad}_{g}\left(i D_{\vartheta}\right)_{X}(x)\right\|_{x}^{2} \quad\left(\boldsymbol{\vartheta} \in \mathbb{R}^{2}\right) .
$$

Furthermore, let us define $\tilde{\Psi}(u, g T, \tau)=\tilde{\Psi}_{m}(u, g T, \tau) \in \mathfrak{t}$ by setting

$$
\tilde{\Psi}(u, g T, \boldsymbol{\vartheta}):=u \operatorname{diag}\left(\operatorname{Ad}_{g^{-1}}\left(\Phi_{G}^{\prime}(m)\right)\right)-\boldsymbol{\nu}, \quad \Phi_{G}^{\prime}(m):=-i \Phi_{G}(m) .
$$

Finally, let us pose

$$
\Psi(u, g T, \boldsymbol{\vartheta}):=\langle\tilde{\Psi}(u, g T), \boldsymbol{\vartheta}\rangle
$$

The following Proposition is proved by a rather lengthy computation, along the lines of those in the proof of Theorem 1.3.3 and in [Pao12].

## Proposition 4.7.1.

$$
\begin{aligned}
i k & {\left[u \psi\left(\tilde{\mu}_{g e^{-i \boldsymbol{\vartheta} / \sqrt{k}} g^{-1}}\left(x_{\tau, k}\right), x_{\tau, k}\right)-\langle\boldsymbol{\vartheta}, \boldsymbol{\nu}\rangle / \sqrt{k}\right] } \\
& =i \sqrt{k} \Psi(u, g T, \boldsymbol{\vartheta})-\frac{u}{2} \boldsymbol{\vartheta}^{t} E(g T) \boldsymbol{\vartheta}+2 i u \tau\left\langle\boldsymbol{\rho}_{g T}, \boldsymbol{\vartheta}\right\rangle+k R_{3}\left(\frac{\tau}{\sqrt{k}}, \frac{\boldsymbol{\vartheta}}{\sqrt{k}}\right) .
\end{aligned}
$$

Corollary 4.7.1. (4.64) may be rewritten as follows:

$$
\begin{align*}
& \Pi_{k \nu}\left(x_{\tau, k}, x_{\tau, k}\right) \sim  \tag{4.65}\\
& \quad \frac{k\left(\nu_{1}-\nu_{2}\right)}{(2 \pi)^{2}} \int_{G / T} \mathrm{dV}_{G / T}(g T) \int_{-\infty}^{+\infty} \mathrm{d} \vartheta_{1} \int_{-\infty}^{+\infty} \mathrm{d} \vartheta_{2} \int_{0}^{+\infty} \mathrm{d} u \\
& \quad\left[e^{i \sqrt{k} \Psi(u, g T, \boldsymbol{\vartheta})} \mathcal{A}_{k, \boldsymbol{\nu}}(u, g T, \tau, \boldsymbol{\vartheta})\right] .
\end{align*}
$$

where (leaving implicit the dependence on $x$ )

$$
\begin{align*}
\mathcal{A}_{k, \boldsymbol{\nu}}(u, g T, \tau, \boldsymbol{\vartheta}):= & e^{-\frac{u}{2} \vartheta^{t} E(g T) \vartheta+2 i u \tau\left\langle\boldsymbol{\rho}_{g T}, \boldsymbol{\vartheta}\right\rangle+k R_{3}\left(\frac{\tau}{\sqrt{k}}, \frac{\vartheta}{\sqrt{k}}\right)} \Delta\left(e^{i \boldsymbol{\vartheta} / \sqrt{k}}\right)  \tag{4.66}\\
& \cdot s\left(\tilde{\mu}_{g e^{-i \vartheta / \sqrt{k}} g^{-1}}\left(x_{\tau, k}\right), x_{\tau, k}, k u\right) .
\end{align*}
$$

Let $h_{m} T \in G / T$ be the unique coset such that $h_{m}^{-1} \Phi_{G}(m) h_{m}$ is diagonal. Then only a rapidly decreasing contribution to the asymptotics is lost in (4.65), if integration in $\mathrm{dV}_{G / T}$ is localized in a small neighbourhood of $h_{m} T$. In the following, a $C^{\infty}$ bump function on $G / T$, supported in a small neighbourhood of $h_{m} T$ and identically equal to 1 near $h_{m} T$, will be implicitly incorporated into the amplitude (4.66).

For some choice of $h_{m} \in h_{m} T$ and $\delta>0$ sufficiently small, let us consider the real-analytic map

$$
h: w \in B(0 ; \delta) \subset \mathbb{C} \mapsto h(w):=h_{m} \exp \left(i\left(\begin{array}{cc}
0 & w \\
\bar{w} & 0
\end{array}\right)\right) \in G .
$$

By composition with the projection $\pi^{\prime}: G \rightarrow G / T$, we obtain a real-analytic coordinate chart on $G / T$ centred at $h_{m} T \in G / T$, given by $w \in B(0 ; \delta) \mapsto$ $h(w) T \in G / T$. The Haar volume form on $G / T$ has the form $\mathcal{V}_{G / T}(w) \mathrm{dV}_{\mathbb{C}}(w)$, where $\mathrm{dV}_{\mathbb{C}}(w)$ is the Lebesgue measure on $\mathbb{C}$, and $\mathcal{V}_{G / T}$ is a uniquely determined $C^{\infty}$ positive function on $B(0 ; \delta)$. We record the following statements, whose proofs we shall omit for the sake of brevity.

Lemma 4.7.1. $\mathcal{V}_{G / T}$ is rotationally invariant, that is,

$$
\mathcal{V}_{G / T}(w)=\mathcal{V}_{G / T}\left(e^{i \theta} w\right)
$$

for all $w \in B(0 ; \delta)$ and $e^{i \theta} \in S^{1}$. In particular, $\mathcal{V}_{G / T}$ is given by a convergent power series in $r^{2}=|w|^{2}$ on $B(0 ; \delta)$.

Thus we shall write

$$
\mathcal{V}_{G / T}(w)=\mathcal{V}_{G / T}(r)=D_{G / T} \cdot \mathcal{S}_{G / T}(r)
$$

where $D_{G / T}$ is a constant, and $\mathcal{S}_{G / T}(r)=1+\sum_{j} s_{j} r^{2 j}$.
Lemma 4.7.2. Let $V_{3}$ be the total area of the unit sphere $S^{3} \subset \mathbb{C}^{2}$. Then

$$
D_{G / T}=2 \pi V_{3}^{-1}
$$

Furthermore, let us introduce the real-analytic function

$$
\kappa=\kappa_{m}: w \in B(0 ; \delta) \mapsto \operatorname{diag}\left(\operatorname{Ad}_{h(w)^{-1}}\left(\Phi_{G}^{\prime}(m)\right)\right) \in \mathbb{R}^{2}
$$

Then we also have the following.
Lemma 4.7.3. $\kappa$ is rotationally invariant, and is given by a convergent power series of the following form

$$
\kappa(w)=\lambda_{\nu}(m)\left[\boldsymbol{\nu}-r^{2}\left(\nu_{1}-\nu_{2}\right) S_{\kappa}(r) \mathbf{b}\right], \quad \mathbf{b}=\binom{1}{-1}
$$

where $r=|w|$ and $S_{\kappa}(r)$ is a real-analytic function of $r$, of the form

$$
S_{\kappa}(r)=1+\sum_{j \geq 1} b_{j} r^{2 j}
$$

If $w=r e^{i \theta}$ in polar coordinates, we shall write accordingly $\mathcal{V}_{G / T}=$ $\mathcal{V}_{G / T}(r)$ and $\kappa=\kappa(r)$.

Recalling the previous notations, let us set

$$
\tilde{\Psi}_{w}(u):=u \kappa(r)-\boldsymbol{\nu}, \quad \Psi_{w}(u, \boldsymbol{\vartheta}):=\left\langle\tilde{\Psi}_{w}(u), \boldsymbol{\vartheta}\right\rangle .
$$

We obtain the following integral formula (dependence on $x$ on the right hand sides is left implicit).

Proposition 4.7.2. As $k \rightarrow+\infty$ we have

$$
\begin{equation*}
\Pi_{k \nu}\left(x_{\tau, k}, x_{\tau, k}\right) \quad \sim D_{G / T} \frac{k\left(\nu_{1}-\nu_{2}\right)}{(2 \pi)^{2}} \int_{-\pi}^{+\pi} \mathrm{d} \theta \int_{0}^{+\infty} \mathrm{d} r\left[I_{k}(\tau, r, \theta)\right] \tag{4.67}
\end{equation*}
$$

where

$$
\begin{align*}
I_{k}(\tau, r, \theta)= & I_{k}(\tau, w):=\int_{-\infty}^{+\infty} \mathrm{d} \vartheta_{1} \int_{-\infty}^{+\infty} \mathrm{d} \vartheta_{2} \int_{0}^{+\infty} \mathrm{d} u  \tag{4.68}\\
& {\left[e^{i \sqrt{k} \Psi_{w}(u, \boldsymbol{\vartheta})} \mathcal{A}_{k, \boldsymbol{\nu}}\left(u, h\left(r e^{i \theta}\right) T, \tau, \boldsymbol{\vartheta}\right) \mathcal{S}_{G / T}(r) r\right] . }
\end{align*}
$$

Our next goal is to produce an asymptotic expansion for $I_{k}(\tau, r, \theta)$.
Definition 4.7.2 Let us set

$$
\mathbf{n}_{1}(r):=\frac{\kappa(r)}{\|\kappa(r)\|}
$$

and let $\mathbf{n}_{2}(r)$ be uniquely determined for $|r|<\delta$ so that $\mathcal{B}_{r}:=\left(\mathbf{n}_{1}(r), \mathbf{n}_{2}(r)\right)$ is a positively oriented orthonormal basis of $\mathbb{R}^{2}$. We shall write the change of basis matrix in the form

$$
M_{\mathcal{C}_{2}}^{\mathcal{B}_{r}}\left(i d_{\mathbb{R}^{2}}\right)=\left(\begin{array}{cc}
C(r) & -S(r)  \tag{4.69}\\
S(r) & C(r)
\end{array}\right),
$$

where $\mathcal{C}_{2}$ is the canonical basis of $\mathbb{R}^{2}$, and denote the change of coordinates by $\boldsymbol{\vartheta}=\zeta_{1} \mathbf{n}_{1}(w)+\zeta_{2} \mathbf{n}_{2}(w)$.

A straightforward computation then yields the following.

Corollary 4.7.2. With $w=r e^{i \theta} \in B(0 ; \delta)$ and $I_{k}(\tau ; w)$ as in (4.68), we have

$$
\begin{equation*}
I_{k}(\tau, w)=\int_{-\infty}^{+\infty} \mathrm{d} \zeta_{2}\left[e^{-i \sqrt{k}\left\langle\boldsymbol{\nu}, \mathbf{n}_{2}(w)\right\rangle \zeta_{2}} J_{k}\left(\tau, w ; \zeta_{2}\right) \mathcal{S}_{G / T}(r) r\right] \tag{4.70}
\end{equation*}
$$

where

$$
\begin{equation*}
J_{k}\left(\tau, w ; \zeta_{2}\right):=\int_{-\infty}^{+\infty} \mathrm{d} \zeta_{1} \int_{0}^{+\infty} \mathrm{d} u\left[e^{i \sqrt{k} \Upsilon_{r}\left(u, \zeta_{1}\right)} \mathcal{A}_{k, \boldsymbol{\nu}}(u, h(w) T, \tau, \boldsymbol{\vartheta}(\boldsymbol{\zeta}))\right] \tag{4.71}
\end{equation*}
$$

and

$$
\Upsilon_{r}\left(u, \zeta_{1}\right):=\left[u\|\kappa(r)\|-\left\langle\boldsymbol{\nu}, \mathbf{n}_{1}(r)\right\rangle\right] \zeta_{1} .
$$

Let us view $J_{k}$ as an oscillatory integral with phase $\Upsilon_{r}$.
Lemma 4.7.4. $\Upsilon_{r}$ has the unique critical point

$$
P_{r}=(u(r), 0):=\left(\frac{\left\langle\boldsymbol{\nu}, \mathbf{n}_{1}(r)\right\rangle}{\|\kappa(r)\|}, 0\right)
$$

Furthermore, $\Upsilon_{r}\left(P_{r}\right)=0$, and the Hessian matrix is

$$
H\left(\Upsilon_{r}\right)_{P_{r}}=\left(\begin{array}{cc}
0 & \|\kappa(r)\| \\
\|\kappa(r)\| & 0
\end{array}\right)
$$

Hence its signature is zero and the critical point is non-degenerate.
Recalling that $s_{0}(x, x)=\pi^{-d}$, the amplitude in (4.71) may be written in the following form

$$
\begin{align*}
& \mathcal{A}_{k, \boldsymbol{\nu}}(u, h(w) T, \tau, \boldsymbol{\vartheta}(\boldsymbol{\zeta}))  \tag{4.72}\\
& \quad \sim e^{-\frac{u}{2} \boldsymbol{\vartheta}(\boldsymbol{\zeta})^{t} E(w) \boldsymbol{\vartheta}(\boldsymbol{\zeta})+2 i u \tau\left\langle\boldsymbol{\rho}_{h(w) T}, \boldsymbol{\vartheta}(\zeta)\right\rangle}\left[e^{\frac{i}{\sqrt{k}} \boldsymbol{\vartheta}_{1}(\boldsymbol{\zeta})}-e^{\frac{i}{\sqrt{k}} \boldsymbol{\vartheta}_{2}(\boldsymbol{\zeta})}\right] \\
& \quad \cdot\left(\frac{k u}{\pi}\right)^{d}\left[1+\sum_{j \geq 1} a_{j}(u, w ; \tau, \boldsymbol{\vartheta}(\boldsymbol{\zeta})) k^{-j / 2}\right]
\end{align*}
$$

in (4.72) we have set $E(w):=\tilde{E}(h(w) T)$, and in view of the exponent $k R_{3}(\tau / \sqrt{k}, \boldsymbol{\vartheta} / \sqrt{k})$ appearing in (4.66), $a_{j}(u, w ; \cdot, \cdot)$ is an appropriate polynomial in $(\tau, \boldsymbol{\vartheta})$ of degree $\leq 3 j$.

Given Lemma 4.7.4, we may evaluate $J_{k}$ in (4.71) by the Stationary Phase Lemma, and obtain an asymptotic expansion in descending powers of $k^{1 / 2}$. The latter expansion may be inserted in (4.70), and integrated term by term, thus leading to an asymptotic expansion for $I_{k}$. The leading order term of either expansion is determined by the contribution of the leading order term in the asymptotic expansion for the amplitude in (4.72), which is given by the following:

$$
\begin{aligned}
& J_{k}^{\prime}\left(\tau, w ; \zeta_{2}\right)=\left(\frac{k}{\pi}\right)^{d} \int_{-\infty}^{+\infty} \mathrm{d} \zeta_{1} \int_{0}^{+\infty} \mathrm{d} u\left[e^{i \sqrt{k} \Upsilon_{w}\left(u, \zeta_{1}\right)} u^{d}\left(e^{\frac{i}{\sqrt{k}} \vartheta_{1}(\zeta)}-e^{\frac{i}{\sqrt{k}} \vartheta_{2}(\zeta)}\right)\right. \\
&\left.\cdot e^{-\frac{u}{2} \vartheta(\zeta)^{t} E(w) \vartheta(\zeta)+2 i u \tau\left\langle\boldsymbol{\rho}_{h(w) T}, \boldsymbol{\vartheta}(\zeta)\right\rangle}\right]
\end{aligned}
$$

Definition 4.7.3. Suppose $w=r e^{i \theta} \in B(0, \delta)$ and let $C(r)$ and $S(r)$ be as in (4.69). Let us set

$$
\begin{aligned}
\mathfrak{a}(w) & :=u(r)\left(\begin{array}{ll}
-S(r) & C(r)) E(w)
\end{array}\binom{-S(r)}{C(r)}\right. \\
& =u(r)\left\|\operatorname{Ad}_{h(w)}\left(\mathbf{n}_{2}(r)\right)_{X}(x)\right\|_{x}^{2}
\end{aligned}
$$

and

$$
\begin{aligned}
\mathfrak{r}(w) & :=2 u(r)\left\langle\boldsymbol{\rho}_{h(w) T}, \mathbf{n}_{2}(r)\right\rangle \\
& =2 u(r) \omega_{m}\left(\operatorname{Ad}_{h(w)}\left(\mathbf{n}_{2}(r)\right)_{M}(m), \Upsilon_{\nu}(m)\right) .
\end{aligned}
$$

Given the previous considerations, an application of the Stationary Phase Lemma yields the following.

Definition 4.7.4. With $|r|<\delta$, let us set $\mathfrak{b}(r):=\left\langle\boldsymbol{\nu}, \mathbf{n}_{2}(r)\right\rangle$, and

$$
D_{l}(r):=\frac{i^{l}}{l!\|\kappa(r)\|}\left[C(r)^{l}+(-1)^{l-1} S(r)^{l}\right]
$$

The definition of $\mathfrak{b}(r)$ implies:

$$
\mathfrak{b}(r)=-\frac{\left(\nu_{1}-\nu_{2}\right)\left(\nu_{1}+\nu_{2}\right)}{\|\boldsymbol{\nu}\|} r^{2} S_{1}(r)
$$

where $S_{1}$ is a real-analytic function of the form $S_{1}(r)=1+\sum_{j \geq 1} c_{j} r^{2 j}$.

Proposition 4.7.3. Suppose $x \in X_{\mathcal{O}_{\nu}}^{G}$. Then as $k \rightarrow+\infty$ we have

$$
\begin{align*}
& \Pi_{k \nu}\left(x_{\tau, k}, x_{\tau, k}\right)  \tag{4.73}\\
& \quad \sim D_{G / T} \frac{k\left(\nu_{1}-\nu_{2}\right)}{(2 \pi)^{2}} \int_{-\pi}^{\pi} \mathrm{d} \theta \int_{0}^{+\infty} \mathrm{d} r\left[I_{k}(\tau, r, \theta)\right]
\end{align*}
$$

where $I_{k}(\tau ; r ; \theta)$ is given by an asymptotic expansion in descending powers of $k^{1 / 2}$, the leading power being $k^{d-1}$. As a function of $\tau$, aside from a phase factor, the coefficient of $k^{d-(1+j) / 2}$ is a polynomial of degree $\leq 3 j$. Up to non-dominant terms we may replace $I_{k}(\tau, w)$ by

$$
\begin{align*}
I_{k}(\tau, w)^{\prime}= & -\left(\frac{k}{\pi}\right)^{d}\left(\frac{2 \pi}{\sqrt{k}}\right) \mathcal{S}_{G / T}(r) r \cdot u(w)^{d}  \tag{4.74}\\
& \cdot \sum_{l \geq 1} \frac{D_{l}(r)}{k^{l / 2}} \int_{-\infty}^{+\infty} d \zeta_{2}\left[e^{-i \sqrt{k} \zeta_{2} \mathfrak{f}(\tau, w)} \zeta_{2}^{l} \cdot e^{-\frac{1}{2} \mathfrak{a}(w) \zeta_{2}^{2}}\right]
\end{align*}
$$

where for $k=1,2, \ldots$, we have set

$$
\begin{equation*}
\mathfrak{f}_{k}(\tau, w):=\mathfrak{b}(r)-\frac{\tau}{k^{1 / 2}} \mathfrak{r}(w) \tag{4.75}
\end{equation*}
$$

The Gaussian integrals in (4.74) may be estimated recalling that

$$
\begin{equation*}
\int_{-\infty}^{+\infty} x^{l} e^{-i \xi x-\frac{1}{2} \lambda x^{2}} \mathrm{~d} x=\sqrt{2 \pi} \frac{(-i)^{l}}{\lambda^{l+1 / 2}} P_{l}(\xi) e^{-\frac{1}{2 \lambda} \xi^{2}} \tag{4.76}
\end{equation*}
$$

where $P_{l}(\xi)=\xi^{l}+\sum_{j \geq 1} p_{l_{j}} \xi^{l-2 j}$ is a monic polynomial in $\xi$, of degree $l$ and parity $(-1)^{l}$ (thus the previous sum is finite). Applying (4.76) with

$$
\xi=k^{1 / 2} \mathfrak{f}_{k}(w, \tau), \quad \lambda=\mathfrak{a}(w)
$$

we obtain the following conclusion.
Proposition 4.7.4. Let us set

$$
\begin{equation*}
F_{l}(\tau, w):=\frac{\sqrt{2 \pi}}{l!}\left[\frac{C(r)^{l}+(-1)^{l-1} S(r)^{l}}{\|\kappa(r)\|}\right] \frac{P_{l}\left(\sqrt{k} \mathfrak{f}_{k}(\tau, w)\right)}{k^{l / 2} \mathfrak{a}(w)^{l+1 / 2}} \tag{4.77}
\end{equation*}
$$

Up to lower order terms, we can replace $I_{k}^{\prime}$ in (4.74) by

$$
\begin{align*}
I_{k}(\tau, w)^{\prime \prime}:= & -\left(\frac{k}{\pi}\right)^{d}\left(\frac{2 \pi}{\sqrt{k}}\right) \mathcal{S}_{G / T}(r) \cdot r \cdot u(w)^{d}  \tag{4.78}\\
& \cdot e^{-\frac{1}{2} \frac{\mathfrak{f}_{k}(\tau, w)^{2}}{\mathfrak{a}(w)}} \sum_{l \geq 1} F_{l}(\tau, w)
\end{align*}
$$

Thus the leading order asymptotics of $\Pi_{k \nu}\left(x_{\tau, k}, x_{\tau, k}\right)$ are obtained by replacing $I_{k}(\tau, r, \theta)$ in (4.73) by $I_{k}(\tau, w)^{\prime \prime}$.

### 4.7.2 Proof of Theorem 1.3.4

We shall set $\tau=0$ in 4.7.3 and obtain an asymptotic estimate for $\Pi_{k \nu}(x, x)$ when $x \in X_{\mathcal{O}_{\nu}}^{G}$ and $k \rightarrow+\infty$.

Proof of Theorem 1.3.4. It follows from the definitions that

$$
\frac{\mathfrak{f}_{k}(0, w)^{2}}{\mathfrak{a}(w)}=\frac{\mathfrak{b}(r)^{2}}{\mathfrak{a}(w)}=\lambda_{\nu}(m) D(\boldsymbol{\nu}) r^{4} \mathcal{S}(r, \theta)
$$

where $\mathcal{S}(r, \theta)=1+\sum_{j \geq 1} r^{j} d_{j}(\theta)$, and

$$
D(\boldsymbol{\nu}):=\frac{\left(\nu_{1}-\nu_{2}\right)^{2}\left(\nu_{1}+\nu_{2}\right)^{2}}{\left\|\operatorname{Ad}_{h_{m}}\left(\boldsymbol{\nu}_{\perp}\right)_{M}(m)\right\|_{m}^{2}}
$$

Similarly,

$$
\begin{aligned}
\frac{P_{l}\left(\sqrt{k} \mathfrak{f}_{k}(0, w)\right)}{k^{l / 2} \mathfrak{a}(w)^{l+1 / 2}} & =\frac{P_{l}(\sqrt{k} \mathfrak{b}(r))}{k^{l / 2} \mathfrak{a}(w)^{l+1 / 2}} \\
& =\frac{1}{\mathfrak{a}(w)^{l+1 / 2}}\left[\mathfrak{b}(r)^{l}+\sum_{j \geq 1}^{\lfloor l / 2\rfloor} p_{l_{j}} k^{-j} \mathfrak{b}(r)^{l-2 j}\right] \\
& =\sum_{j=0}^{\lfloor l / 2\rfloor} \frac{1}{k^{j}} r^{2 l-4 j} \mathcal{S}_{l_{j}}(r, \theta)
\end{aligned}
$$

where $\mathcal{S}_{l_{j}}(r, \theta)$ is a convergent power series in $r$. The resulting series may be integrated term by term. The $l$-th summand in (4.78) then gives rise to a convergent series of summands of the form

$$
B_{\boldsymbol{\nu}, l, j}(m, \theta) \frac{1}{k^{j}} \int_{0}^{+\infty} \tilde{r}^{2 l-4 j+a} e^{-\frac{1}{2} k \lambda_{\boldsymbol{\nu}}(m) D(\boldsymbol{\nu}) \cdot \tilde{r}^{4}} \tilde{r} \mathrm{~d} \tilde{r}=O\left(\frac{1}{k^{\frac{l+1}{2}+\frac{a}{4}}}\right)
$$

with $j \leq\lfloor l / 2\rfloor$ and $a=0,1,2, \ldots$.
The previous discussion shows that $\Pi_{k \nu}(x, x)$ is given by an asymptotic expansion in descending powers of $k^{1 / 4}$, and that the leading order term occurs for $l=1$ and $a=0$.

By equation (4.76), $P_{l}(\xi)=\xi$; by Lemma 4.7.3, $\|\kappa(r)\|=\lambda_{\boldsymbol{\nu}}(m)\|\boldsymbol{\nu}\| \cdot$ $\mathcal{S}_{\kappa}^{\prime}(r)$, where $\mathcal{S}_{\kappa}^{\prime}(r)$ is a convergent power series in $r^{2}$ with $\mathcal{S}_{\kappa}^{\prime}(0)=1$.

In view of (4.77), we obtain

$$
F_{l}(0, w)=-\sqrt{2 \pi} \cdot \frac{\left(\nu_{1}-\nu_{2}\right)\left(\nu_{1}+\nu_{2}\right)^{2}}{\left\|\operatorname{Ad}_{h_{m}}\left(\boldsymbol{\nu}_{\perp}\right)_{M}(m)\right\|^{3}} \lambda_{\nu}(m)^{1 / 2} r^{2} \mathcal{S}_{F_{1}}(r, \theta)
$$

where $\mathcal{S}_{F_{1}}$ is real-analytic function and $\mathcal{S}^{\prime \prime}(0, \theta)=1$.
Hence the leading order term of the asymptotic expansion of $\Pi_{k \nu}(x, x)$ is given by

$$
\begin{equation*}
D_{G / T} \frac{k\left(\nu_{1}-\nu_{2}\right)}{(2 \pi)^{2}} \int_{-\pi}^{\pi} \mathrm{d} \theta \int_{0}^{+\infty} \mathrm{d} r\left[L_{k}(r, \theta)\right] \tag{4.79}
\end{equation*}
$$

where

$$
\begin{align*}
L_{k}(r):= & 2^{3 / 2} \frac{k^{d-1 / 2}}{\pi^{d-3 / 2}} \lambda_{\nu}(m)^{-(d-1 / 2)}  \tag{4.80}\\
& \cdot\left[\frac{\left(\nu_{1}-\nu_{2}\right)\left(\nu_{1}+\nu_{2}\right)^{2}}{\left\|\operatorname{Ad}_{h_{m}}\left(\boldsymbol{\nu}_{\perp}\right)_{M}(m)\right\|^{3}}\right] e^{-\frac{1}{2} k \lambda_{\nu}(m) D(\boldsymbol{\nu}) r^{4} \mathcal{S}(r, \theta)} r^{3} \tilde{\mathcal{S}}(r, \theta),
\end{align*}
$$

where again $\tilde{\mathcal{S}}$ is real-analytic and e $\tilde{\mathcal{S}}(0, \theta) \equiv 1$.
We need to integrate in $\mathrm{d} r$ the product of the last two factors in (4.80). Let us perform the coordinate change $s=\sqrt{k} r^{2} \mathcal{S}(r, \theta)^{1 / 2}$, and argue as above. To leading order, we are reduced to computing

$$
\frac{1}{2 k} \int_{0}^{+\infty} \mathrm{d} s\left[e^{-\frac{1}{2} \lambda_{\boldsymbol{\nu}}(m) D(\boldsymbol{\nu}) s^{2}} s\right]=\frac{1}{2 k} \cdot \frac{1}{\lambda_{\boldsymbol{\nu}}(m) D(\boldsymbol{\nu})}
$$

Inserting this in (4.79), we conclude that the leading order term in the asymptotic expansion of $\Pi_{k \nu}(x, x)$ is

$$
\frac{D_{G / T}}{\sqrt{2}} \frac{1}{\|\Phi(m)\|^{d+1 / 2}}\left(\frac{k\|\boldsymbol{\nu}\|}{\pi}\right)^{d-1 / 2} \cdot \frac{\|\boldsymbol{\nu}\|}{\left\|\operatorname{Ad}_{h_{m}}\left(\boldsymbol{\nu}_{\perp}\right)_{M}(m)\right\|}
$$

The proof of Theorem 1.3.4 is complete.

### 4.7.3 Proof of Theorem 1.3.6

The proof is a modification of the one of Theorem 1.3.4, so the discussion will be sketchy. We shall set

$$
x_{j, k}:=x+\frac{1}{\sqrt{k}} \mathbf{v}_{j}, \quad j=1,2
$$

Definition 4.7.5. With the previous notation, let us set

$$
\begin{aligned}
\Gamma(\boldsymbol{\vartheta}, & \left.g T, \mathbf{v}_{j}\right) \\
:= & -\frac{1}{2}\left[\left\langle\operatorname{diag}\left(\operatorname{Ad}_{g^{-1}}\left(\Phi_{G}^{\prime}(m)\right)\right), \boldsymbol{\vartheta}\right\rangle^{2}+\left\|\mathbf{v}_{1}-\mathbf{v}_{2}+\operatorname{Ad}_{g}\left(i D_{\vartheta}\right)_{M}(m)\right\|_{m}^{2}\right] \\
& +i\left[-\omega_{m}\left(\mathbf{v}_{1}, \mathbf{v}_{2}\right)+\omega_{m}\left(\operatorname{Ad}_{g}\left(i D_{\vartheta}\right)_{M}(m), \mathbf{v}_{1}+\mathbf{v}_{2}\right)\right]
\end{aligned}
$$

Then, the same computations leading to Proposition 4.7.1 yield the following.

## Proposition 4.7.5.

$$
\begin{aligned}
i k & {\left[u \psi\left(\tilde{\mu}_{g e^{-i \vartheta / \sqrt{k}} g^{-1}}\left(x_{1, k}\right), x_{2, k}\right)-\langle\boldsymbol{\vartheta}, \boldsymbol{\nu}\rangle / \sqrt{k}\right] } \\
& =i \sqrt{k} \Psi(u, g T, \boldsymbol{\vartheta})+u \Gamma\left(\boldsymbol{\vartheta}, g T, \mathbf{v}_{j}\right)+k R_{3}\left(\frac{\mathbf{v}_{j}}{\sqrt{k}}, \frac{\boldsymbol{\vartheta}}{\sqrt{k}}\right) .
\end{aligned}
$$

Remark 5. Assuming $\mathbf{v}_{1}, \mathbf{v}_{2} \in \mathfrak{g}_{M}\left(m_{x}\right)^{\perp_{h}}$ and recalling Definition 4.7.1 we have

$$
\Gamma\left(\boldsymbol{\vartheta}, g T, \mathbf{v}_{j}\right)=\psi_{2}\left(\mathbf{v}_{1}, \mathbf{v}_{2}\right)-\frac{1}{2} \boldsymbol{\vartheta}^{t} E(g T) \boldsymbol{\vartheta}
$$

In place of Corollary 4.7.1, we then obtain the following:

$$
\begin{aligned}
& \Pi_{k \nu}\left(x_{1, k}, x_{2, k}\right) \sim \\
& \frac{k\left(\nu_{1}-\nu_{2}\right)}{(2 \pi)^{2}} \int_{G / T} \mathrm{dV}_{G / T}(g T) \int_{-\infty}^{+\infty} \mathrm{d} \vartheta_{1} \int_{-\infty}^{+\infty} \mathrm{d} \vartheta_{2} \int_{0}^{+\infty} \mathrm{d} u \\
& \quad\left[e^{i \sqrt{k} \Psi(u, g T, \boldsymbol{\vartheta})} \mathcal{A}_{k, \boldsymbol{\nu}}^{\prime}\left(u, g T, \boldsymbol{\vartheta}, \mathbf{v}_{j}\right)\right] .
\end{aligned}
$$

with the amplitude

$$
\begin{aligned}
\mathcal{A}_{k, \boldsymbol{\nu}}^{\prime}\left(u, g T, \boldsymbol{\vartheta}, \mathbf{v}_{j}\right):= & e^{u \psi_{2}\left(\mathbf{v}_{1}, \mathbf{v}_{2}\right)-\frac{u}{2} \vartheta^{t} E(g T) \vartheta+k R_{3}\left(\frac{\mathbf{v}_{j}}{\sqrt{k}} \frac{\vartheta}{\sqrt{k}}\right)} \Delta\left(e^{i \vartheta / \sqrt{k}}\right) \\
& \cdot s\left(\tilde{\mu}_{g e^{-i \vartheta / \sqrt{k}} g^{-1}}\left(x_{1, k}\right), x_{2, k}, k u\right)
\end{aligned}
$$

Similarly in place of (4.72) we now have the following expansion:

$$
\begin{aligned}
& \mathcal{A}_{k, \boldsymbol{\nu}}^{\prime}\left(u, g T, \boldsymbol{\vartheta}, \mathbf{v}_{j}\right) \\
& \quad \sim e^{u \psi_{2}\left(\mathbf{v}_{1}, \mathbf{v}_{2}\right)-\frac{u}{2} \boldsymbol{\vartheta}^{t} E(g T) \vartheta+k R_{3}\left(\frac{\mathbf{v}_{j}}{\sqrt{k}} \frac{\vartheta}{\sqrt{k}}\right)}\left[e^{\frac{i}{\sqrt{k}} \boldsymbol{\vartheta}_{1}(\boldsymbol{\zeta})}-e^{\frac{i}{\sqrt{k}} \boldsymbol{\vartheta}_{2}(\boldsymbol{\zeta})}\right] \\
& \quad \cdot\left(\frac{k u}{\pi}\right)^{d} \cdot\left[1+\sum_{j \geq 1} a_{j}\left(u, w ; \mathbf{v}_{1}, \mathbf{v}_{2}, \boldsymbol{\vartheta}(\boldsymbol{\zeta})\right) k^{-j / 2}\right]
\end{aligned}
$$

where $a_{j}$ is, as a function of $\mathbf{v}_{1}$ and $\mathbf{v}_{2}$, a polynomial of degree $\leq 3 j$.
With these changes, Theorem 1.3.6 can be proved by applying the arguments in the proof of Theorem 1.3.4 with minor modifications.

### 4.7.4 Proof of Theorem 1.3.7

Proof. Let $A^{\prime} \subset X$ be a one-sided "outer" tubular neighbourhood of $X_{\mathcal{O}_{\nu}}^{G}$, that is, the intersection of $A$ with a tubular neighbourhood of $X_{\mathcal{O}_{\nu}}^{G}$ in $X$.

By Theorem 1.3.1, we have

$$
\begin{aligned}
& \operatorname{dim}_{\text {out }} H(X)_{k \nu} \\
& \quad=\int_{A} \Pi_{k \nu}(x, x) \mathrm{d} V_{X}(x) \sim \int_{A^{\prime}} \Pi_{k \nu}(x, x) \mathrm{dV}_{X}(x) .
\end{aligned}
$$

Let us denote by $m_{x}=\pi(x)$ the sign of $\nu_{1}+\nu_{2}$. Then, locally along $X_{\mathcal{O}_{\nu}}^{G}$, for some sufficiently small $\delta>0$ we can parametrize $A^{\prime}$ by a diffeomorphism

$$
\Gamma: X_{\mathcal{O}_{\nu}}^{G} \times[0, \delta) \rightarrow A^{\prime}, \quad(x, \tau) \mapsto x+\tau(\boldsymbol{\nu}) \Upsilon_{\nu}\left(m_{x}\right)
$$

where $m_{x}=\pi(x)$. The latter expression is meant in terms of a collection of smoothly varying systems of Heisenberg local coordinates centred at $x \in$ $X_{\mathcal{O}_{\nu}}^{G}$, locally defined along $X_{\mathcal{O}_{\nu}}^{G}$ (to be precise, one ought to work locally on $X_{\mathcal{O}_{\nu}}^{G}$, introduce an appropriate open cover of $X_{\mathcal{O}_{\nu}}^{G}$, and a subordinate partition of unity; however for the sake of exposition we shall omit details on this).

We shall set $x_{\tau}:=\Gamma(x, \tau)$, and write

$$
\Gamma^{*}\left(\mathrm{~d} V_{X}\right)=\mathcal{V}_{X}(x, \tau) \mathrm{d} V_{X_{\mathcal{O}_{\nu}}}(x) \mathrm{d} \tau
$$

where $\mathcal{V}_{X}: X_{\mathcal{O}_{\nu}}^{G} \times[0, \delta) \rightarrow(0,+\infty)$ is $\mathcal{C}^{\infty}$ and $\mathcal{V}_{X}(x, 0)=\left\|\Upsilon_{\nu}\left(m_{x}\right)\right\|$.
Hence we obtain

$$
\begin{align*}
& \operatorname{dim}_{\text {out }} H(X)_{k \nu}  \tag{4.81}\\
& \quad \sim \int_{X_{O_{\nu}}} \mathrm{dV} V_{X_{\boldsymbol{O}}^{G}}(x) \int_{0}^{\delta} \mathrm{d} \tau\left[\mathcal{V}_{X}(x, \tau) \Pi_{k \nu}\left(x_{\tau}, x_{\tau}\right)\right]
\end{align*}
$$

By Theorem 1.3.3, only a rapidly decreasing contribution to (4.81) is lost, if integration in (4.81) is restricted to the locus where $\tau \leq C k^{\epsilon-1 / 2}$. Thus the asymptotics of $\operatorname{dim}_{o u t} H(X)_{k \nu}$ are unchanged, if the integrand is multiplied
by a rescaled cut-off function $\varrho\left(k^{1 / 2-\epsilon} \tau\right)$, where $\varrho$ is identically one sufficiently near the origin in $\mathbb{R}$, and vanishes outside a slightly larger neighbourhood.

With the rescaling $\tau \mapsto \tau / \sqrt{k}$, we obtain

$$
\operatorname{dim}_{\text {out }} H(X)_{k \nu} \sim \frac{1}{\sqrt{k}} \int_{X_{G_{\nu}}^{G}} \mathrm{dV}_{X_{\mathcal{O}_{\nu}}}(x)\left[\mathcal{H}_{k}(x)\right]
$$

where with $x_{\tau, k}:=\Gamma\left(x, k^{-1 / 2} \tau\right)$ we have set

$$
\begin{equation*}
\mathcal{H}_{k}(x):=\int_{0}^{+\infty} \mathrm{d} \tau\left[\varrho\left(k^{-\epsilon} \tau\right) \mathcal{V}_{X}\left(x, \frac{\tau}{\sqrt{k}}\right) \Pi_{k \nu}\left(x_{\tau, k}, x_{\tau, k}\right)\right] . \tag{4.82}
\end{equation*}
$$

Integration in $\mathrm{d} \tau$ is now over an expanding interval of the form $\left[0, C^{\prime} k^{\epsilon}\right)$.
Let us consider the asymptotics of 4.82. Having in mind (4.78), and inserting the Taylor expansion of $\mathcal{V}_{X}$, we are led to considering double integrals of the form

$$
\begin{align*}
& \frac{1}{k^{(l+j) / 2}} \int_{0}^{+\infty} \mathrm{d} \tau \int_{0}^{+\infty} \mathrm{d} r  \tag{4.83}\\
& \quad\left[r C(r)^{l} \tau^{j} \mathcal{S}^{\prime}(r) \frac{P_{l}\left(\sqrt{k} \mathfrak{f}_{k}(\tau, w)\right)}{\mathfrak{a}(w)^{l+1 / 2}} \cdot e^{-\frac{1}{2} k \frac{\mathfrak{f}_{k}(\tau, w)^{2}}{\mathfrak{a}(w)}}\right]
\end{align*}
$$

with $l \geq 1$ and $j \geq 0$, and their analogues with $S(r)$ in place of $C(r) ; \mathcal{S}^{\prime}$ is some real-analytic function (dependence on $\theta$ and $x$ is implicit).

In view of (4.75), we have

$$
\frac{\mathfrak{f}_{k}(\sigma(\boldsymbol{\nu}) \tau, w)}{\sqrt{\mathfrak{a}(w)}}=-\sigma(\boldsymbol{\nu})\left[\frac{\left(\nu_{1}-\nu_{2}\right)\left|\nu_{1}+\nu_{2}\right|}{\|\boldsymbol{\nu}\| \sqrt{\mathfrak{a}(0)}} r^{2} S_{1}(r)+\frac{\tau}{k^{1 / 2}} \frac{\mathfrak{r}(0)}{\sqrt{\mathfrak{a}(0)}} S_{2}(r, \theta)\right]
$$

where again $S_{2}(0, \theta)=1$. Therefore, with the change of variables

$$
s:=k^{1 / 4} r \sqrt{S_{1}(r)}, \quad \tilde{\tau}:=\tau S_{2}(r, \theta)
$$

we obtain

$$
\frac{\mathfrak{f}_{k}(\sigma(\boldsymbol{\nu}) \tau, w)}{\sqrt{\mathfrak{a}(w)}}=-\frac{\sigma(\boldsymbol{\nu})}{\sqrt{k}}\left[\frac{\left(\nu_{1}-\nu_{2}\right)\left|\nu_{1}+\nu_{2}\right|}{\|\boldsymbol{\nu}\| \sqrt{\mathfrak{a}(0)}} s^{2}+\frac{\mathfrak{r}(0)}{\sqrt{\mathfrak{a}(0)}} \tilde{\tau}\right] .
$$

Therefore, we also have

$$
\mathfrak{f}_{k}(\sigma(\boldsymbol{\nu}) \tau, w)=-\frac{\sigma(\boldsymbol{\nu})}{\sqrt{k}}\left[\frac{\left(\nu_{1}-\nu_{2}\right)\left|\nu_{1}+\nu_{2}\right|}{\|\boldsymbol{\nu}\|} s^{2}+\mathfrak{r}(0) \tilde{\tau}\right] \cdot\left[1+R_{1}\left(\frac{s}{\sqrt[4]{k}}\right)\right] .
$$

With the substitution $a=s^{2}$, (4.83) may be rewritten as a linear combination of summands of the form

$$
\begin{aligned}
& \frac{1}{k^{(l+j+1) / 2}} \int_{0}^{+\infty} \mathrm{d} \tau \int_{0}^{+\infty} \mathrm{d} a \\
& \quad\left[C\left(\frac{\sqrt{a}}{\sqrt[4]{k}}\right)^{l}\left(A_{1} a+B_{1} \tau\right)^{b} \tau^{j} \cdot\left[1+R_{1}\left(\frac{\sqrt{a}}{\sqrt[4]{k}}\right)\right] \cdot e^{-\frac{1}{2}\left(A_{1} a+B_{1} \tau\right)^{2}}\right] \\
& =O\left(\frac{1}{k^{(l+j+1) / 2}}\right) .
\end{aligned}
$$

Hence the leading contribution occurs for $l=1, j=0$, and dropping the term $R_{1}\left(k^{-1 / 4} \sqrt{a}\right)$. The conclusion of Theorem 1.3.7 then follows by a fairly simple computation.

### 4.8 Proof of Theorem 1.3.5

Arguing as in the proof of Theorem 1.3.4, we can restrict the integration in $\mathrm{dV} V_{G}$ over a small neighbourhood of each $g_{j} \in G_{x}$. Let us denote with $\varrho_{j}$ the bump function vanishing outside a small neighborhood of $g_{j}$ in $G$. We have,

$$
\begin{align*}
& \Pi_{k \nu}(x, x)  \tag{4.84}\\
& \quad \sim \sum_{j=1}^{\left|G_{x}\right|} \frac{k^{2}\left(\nu_{1}-\nu_{2}\right)}{(2 \pi)^{2}} \int_{-\pi}^{\pi} \int_{-\pi}^{\pi} \int_{G / T} \int_{0}^{+\infty} e^{i k \Psi_{x, x}} \mathcal{A}_{j} \mathrm{~d} u \mathrm{~d} V_{G / T}(g T) \mathrm{d} \boldsymbol{\vartheta} ;
\end{align*}
$$

where

$$
\begin{align*}
\Psi_{x, x} & =\Psi_{x, x}(u, \boldsymbol{\vartheta}, g T):=u \psi\left(\tilde{\mu}_{g e^{-i \vartheta} g^{-1}}(x), x\right)-\langle\boldsymbol{\nu}, \boldsymbol{\vartheta}\rangle  \tag{4.85}\\
\mathcal{A}_{j} & =\mathcal{A}_{j}(u, \boldsymbol{\vartheta}, g T):=\Delta\left(e^{i \boldsymbol{\vartheta}}\right) s^{\prime}\left(\tilde{\mu}_{g e^{-i \vartheta} g^{-1}}(x), x, k u\right)
\end{align*}
$$

with

$$
s^{\prime}\left(\tilde{\mu}_{g e^{-i} \vartheta_{g^{-1}}}(x), x, k u\right):=s\left(\tilde{\mu}_{g e^{-i} \vartheta_{g^{-1}}}(x), x, k u\right) \varrho_{j}\left(g e^{-i \vartheta} g^{-1}\right) .
$$

Suppose that $\varrho \in \mathcal{C}_{0}^{\infty}((0,+\infty))$ is $\equiv 1$ on $(1 / D, D)$ and is supported on $(1 /(2 D), 2 D)$. Lemma 4.6.3 holds with minor modifications. Thus the asymptotics of (4.84) is unaltered, if the integrand is multiplied by $\varrho(u)$. We obtain

$$
\begin{align*}
& \Pi_{k \nu}(x, y) \sim  \tag{4.86}\\
& \quad \sum_{j=1}^{\left|G_{x}\right|} \frac{k^{2}\left(\nu_{1}-\nu_{2}\right)}{(2 \pi)^{2}} \int_{-\pi}^{\pi} \int_{-\pi}^{\pi} \int_{G / T} \int_{1 / 2 D}^{2 D} e^{i k \Psi_{x, y}} \mathcal{A}_{j}^{\prime} \mathrm{d} u \mathrm{dV}_{G / T}(g T) \mathrm{d} \boldsymbol{\vartheta}
\end{align*}
$$

where we have set

$$
\mathcal{A}_{j}^{\prime}(u, \boldsymbol{\vartheta}, g T):=\varrho(u) \mathcal{A}_{j}(u, \boldsymbol{\vartheta}, g T)
$$

and the integration in $\mathrm{d} u$ is now over a compact interval. Let us write $\Pi_{k \nu}(x, x)_{j}$ for the $j$-th summand in (4.86).

We shall distinguish two cases, depending on whether $g_{j} \in Z_{x}$ or not. Let us write

$$
\Pi_{k \nu}(x, x)=\Pi_{k \nu}(x, x)_{Z_{x}}+\Pi_{k \nu}(x, x)_{G_{x} \backslash Z_{x}}
$$

where

$$
\Pi_{k \nu}(x, x)_{Z_{x}}:=\sum_{g_{j} \in Z_{x}} \Pi_{k \nu}(x, x)_{j}, \quad \Pi_{k \nu}(x, x)_{G_{x} \backslash Z_{x}}:=\sum_{g_{j} \notin Z_{x}} \Pi_{k \nu}(x, x)_{j}
$$

First, let us consider the summand in (4.86) with $g_{j}$ lying in the centre of the group. Thus, in this case the integral is over a compact neighbourhood of $e^{i \beta_{j}} I_{2}$. Let us perform the following change $\vartheta_{i} \mapsto \vartheta_{i}+\beta_{j}$, for $i=1,2$. We have

$$
\Psi_{x, x}\left(u, \boldsymbol{\vartheta}+\boldsymbol{\beta}_{j}, g T\right)=\Psi_{x, x}(u, \boldsymbol{\vartheta}, g T)-\left(\nu_{1}+\nu_{2}\right) \beta_{j},
$$

where $\boldsymbol{\beta}_{j}^{\boldsymbol{t}}=\left(\beta_{j}, \beta_{j}\right)$. Thus, the $j$-th summand in (4.86) becomes

$$
\begin{align*}
& \Pi_{k \nu}(x, x)_{j} \sim e^{i \beta_{j}\left[1-\left(\nu_{1}+\nu_{2}\right) k\right]} \frac{k^{2}\left(\nu_{1}-\nu_{2}\right)}{(2 \pi)^{2}}  \tag{4.87}\\
& \cdot \int_{-\pi}^{\pi} \mathrm{d} \vartheta_{1} \int_{-\pi}^{\pi} \mathrm{d} \vartheta_{2} \int_{G / T} \mathrm{~d} V_{G / T}(g T) \int_{0}^{+\infty} \mathrm{d} u\left[e^{i k \Psi_{x, x}(u, \boldsymbol{\vartheta}, g T)} \mathcal{A}_{0}^{\prime}(u, \boldsymbol{\vartheta}, g T)\right]
\end{align*}
$$

where, for $j=0$ we have set $g_{0}=I_{2}$. Hence the integral (4.87) is the same appearing in (4.64) with $\tau=0$. Thus we can argue as in the proof of Theorem 1.3.4.

Let us now consider the case $g_{j} \notin Z$. By the discussion preceding the statement of Theorem 1.3.5 in $\S 1.3$, there exist $e^{i \alpha_{j}}, e^{i \beta_{j}} \in S^{1}$ and two unique cosets $h_{m_{x}}, k_{m_{x}} \in G / T$ such that

$$
g_{j}=h_{m_{x}} t_{j 1} h_{m_{x}}^{-1}=k_{m_{x}} t_{j 2} k_{m_{x}}^{-1}
$$

where

$$
k_{m}=h_{m}\left(\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right), \quad t_{j 1}=\left(\begin{array}{cc}
e^{i \alpha_{j}} & 0 \\
0 & e^{i \beta_{j}}
\end{array}\right) \quad \text { and } \quad t_{j 2}=\left(\begin{array}{cc}
e^{i \beta_{j}} & 0 \\
0 & e^{i \alpha_{j}}
\end{array}\right) .
$$

The pulled-back cut-off $\left(g T, e^{i \vartheta}\right) \mapsto \varrho_{j}\left(g e^{\vartheta} g^{-1}\right)$ is supported in two small open neighbourhoods of $\left(h_{m_{x}} T, t_{j 1}\right)$ and $\left(k_{m_{x}} T, t_{j 2}\right)$. Therefore, we can further split $\Pi_{k \nu}(x, x)_{j}, g_{j} \in G_{x} \backslash Z_{x}$,

$$
\Pi_{k \nu}(x, x)_{j}=\Pi_{k \nu}(x, x)_{j 1}+\Pi_{k \nu}(x, x)_{j 2}
$$

where in $\Pi_{k \nu}(x, x)_{j 1}$ (respectively $\left.\Pi_{k \nu}(x, x)_{j 2}\right)$ integration is over a small neighbourhood of $\left(h_{m_{x}} T, t_{j 1}\right)$ (respectively $\left(k_{m_{x}} T, t_{j 2}\right)$ ).

Let us focus on $\Pi_{k \nu}(x, x)_{j 1}$. By the previous remark and by composition with the projection $\pi^{\prime}: G \rightarrow G / T$, we obtain a real-analytic coordinate chart $h(z)$ on $G / T$ centred at $h_{m_{x}} T \in G / T$, and $t\left(\boldsymbol{\vartheta}_{j}\right)$ on $T$ centred at $\left(e^{i \alpha_{j}}, e^{i \beta_{j}}\right)$; explicitly given by the following.

Definition 4.8.1. Let us denote by

$$
h(z):=h_{m_{x}} e^{A(z)}, \quad A(z):=i\left(\begin{array}{cc}
0 & z \\
\bar{z} & 0
\end{array}\right)
$$

and

$$
t\left(\boldsymbol{\vartheta}_{j}\right):=\left(\begin{array}{cc}
e^{i\left(\alpha_{j}+\vartheta_{1}\right)} & 0 \\
0 & e^{i\left(\beta_{j}+\vartheta_{2}\right)}
\end{array}\right)
$$

where $\alpha_{j} \neq \beta_{j}$.
The Haar volume form on $G / T$ has the form $\mathcal{V}_{G / T}(w) \mathrm{dV}_{\mathbb{C}}(w)$, where $\mathrm{d} V_{\mathbb{C}}(w)$ is the Lebesgue measure on $\mathbb{C}$, and $\mathcal{V}_{G / T}$ is a uniquely determined $C^{\infty}$ positive function on $B(0 ; \delta)$. In view of Lemma 4.7.2, we obtain

$$
\mathcal{V}_{G / T}(w)=D_{G / T} \cdot \mathcal{S}_{G / T}(r)
$$

where $D_{G / T}$ is a constant, and $\mathcal{S}_{G / T}(r)=1+\sum_{j} s_{j} r^{2 j}$.

Thus, in this case, the $j$-th summand in (4.86) becomes

$$
\begin{align*}
\Pi_{k \nu}(x, x)_{j 1} \sim & e^{-i\left(\alpha_{j} \nu_{1}+\beta_{j} \nu_{2}\right) k} \cdot \frac{k^{2}\left(\nu_{1}-\nu_{2}\right)}{(2 \pi)^{2}}  \tag{4.88}\\
& \cdot \int_{-\pi}^{\pi} \mathrm{d} \vartheta_{1} \int_{-\pi}^{\pi} \mathrm{d} \vartheta_{2} \int_{\mathbb{C}} \mathrm{d} V_{\mathbb{C}}(z) \int_{0}^{+\infty} \mathrm{d} u \\
& {\left[e^{i k \Psi_{x, x}(u, \boldsymbol{\vartheta}, g T(z))} \mathcal{A}_{j}^{\prime}(u, \boldsymbol{\vartheta}, g T(z)) \mathcal{V}_{G / T}(z)\right] . }
\end{align*}
$$

where

$$
\begin{aligned}
\Psi_{x, x}(u, \boldsymbol{\vartheta}, g T(z)) & :=u \psi\left(\tilde{\mu}_{h(z) t\left(\boldsymbol{\vartheta}_{j}\right)^{-1} h(z)^{-1}}(x), x\right)-\langle\boldsymbol{\nu}, \boldsymbol{\vartheta}\rangle, \\
\mathcal{A}_{j}^{\prime}(u, \boldsymbol{\vartheta}, g T(z)) & :=\Delta\left(t\left(\boldsymbol{\vartheta}_{j}\right)\right) s^{\prime}\left(\tilde{\mu}_{h(z) t\left(\boldsymbol{\vartheta}_{j}\right)^{-1} h(z)^{-1}}(x), x, k u\right) \varrho(u) .
\end{aligned}
$$

As usual let us express $\tilde{\mu}_{h(z) t\left(\boldsymbol{\vartheta}_{j}\right)^{-1} h(z)^{-1}}(x)$ in Heisenberg local coordinates centred in $x$. We have

$$
h(z) t\left(\boldsymbol{\vartheta}_{j}\right)^{-1} h(z)^{-1}=C_{h_{m}}\left[e^{A(z)} e^{-i \boldsymbol{\vartheta}} C_{t_{j 1}^{-1}}\left[e^{-A(z)}\right]\right] g_{j}^{-1},
$$

where $C_{g}$ denotes the conjugation by $g \in G$. Thus, by the Baker-CampbellHausdorff formula, we get

$$
e^{A(z)} e^{-i \vartheta} C_{t_{j 1}^{-1}}\left[e^{-A(z)}\right]=e^{A(z)} e^{-i \vartheta} e^{-\operatorname{Ad}_{t_{j 1}^{-1}} A(z)}=e^{-\gamma_{j}(z, \vartheta)},
$$

where we have set

$$
\gamma_{j}(z, \boldsymbol{\vartheta}):=\gamma_{j 1}(z, \boldsymbol{\vartheta})+\gamma_{j 2}(z, \boldsymbol{\vartheta})+R_{3}(z, \boldsymbol{\vartheta}),
$$

with

$$
\begin{gather*}
\gamma_{j 1}(z, \boldsymbol{\vartheta}):=-A(z)+i \boldsymbol{\vartheta}+\operatorname{Ad}_{t_{j 1}^{-1}} A(z), \\
\gamma_{j 2}(\boldsymbol{\vartheta}, z):=\frac{1}{2}\left[A(z), \operatorname{Ad}_{t_{j 1}^{-1}} A(z)\right]-\frac{i}{2}\left[\boldsymbol{\vartheta}, A(z)+\operatorname{Ad}_{t_{j 1}^{-1}} A(z)\right] \tag{4.89}
\end{gather*}
$$

and $R_{3}(z, \boldsymbol{\vartheta})$ denote a smooth function vanishing at 3-rd order at the origin in $\mathbb{C} \times \mathbb{R}^{2}$. Thus, by Corollary 2.2 of [Pao12], we obtain

$$
\tilde{\mu}_{h(z) t\left(\boldsymbol{\vartheta}_{j}\right)^{-1} h(z)^{-1}}(x)=x+(\Theta(z, \boldsymbol{\vartheta}), V(z, \boldsymbol{\vartheta})),
$$

where

$$
\begin{aligned}
\Theta(z, \boldsymbol{\vartheta}) & :=\left\langle\operatorname{Ad}_{h_{m_{x}^{-1}}^{-1}}\left(\Phi\left(m_{x}\right)\right), \gamma_{j}(z, \boldsymbol{\vartheta})\right\rangle+R_{3}(z, \boldsymbol{\vartheta}), \\
V(z, \boldsymbol{\vartheta}) & :=-\left(\operatorname{Ad}_{h_{m_{x}}} \gamma_{j}(z, \boldsymbol{\vartheta})\right)_{M}\left(m_{x}\right)+R_{2}(z, \boldsymbol{\vartheta}) .
\end{aligned}
$$

In view of the discussion in chapter $\S 3$ of [SZ02], after some computations, we can rewritten the phase $\psi$ as follow

$$
\begin{aligned}
& u \psi\left(\tilde{\mu}_{h(z) t\left(\boldsymbol{\vartheta}_{j}\right)^{-1} h(z)^{-1}}(x), x\right)=u \Theta(z, \boldsymbol{\vartheta}) \\
& \quad+\frac{i u}{2} \cdot\left[\Theta(z, \boldsymbol{\vartheta})^{2}+\|V(z, \boldsymbol{\vartheta})\|^{2}\right]+u R_{3}(z, \boldsymbol{\vartheta}) .
\end{aligned}
$$

Let us notice that $\left(\operatorname{Ad}_{t_{j 1}^{-1}}-I_{2}\right) A(z)$ and the second summand in (4.89) vanish on the diagonal. Thus the phase of the compact supported oscillatory integral $\Pi_{k \nu}(x, x)_{j 1}$ is

$$
\begin{align*}
& \Psi_{x, x}(u, \boldsymbol{\vartheta}, g T)=u\left\langle\operatorname{Ad}_{h(w)^{-1}}(\Phi(m)), \boldsymbol{\vartheta}\right\rangle-\langle\boldsymbol{\nu}, \boldsymbol{\vartheta}\rangle  \tag{4.90}\\
& \quad \frac{i u}{2}\left\|\operatorname{Ad}_{h_{m_{x}}}\left(\gamma_{j}(z, \boldsymbol{\vartheta})\right)_{X}(x)\right\|_{x}^{2}+u\left\langle\operatorname{Ad}_{h(w)^{-1}}(\Phi(m)), \frac{1}{2}\left[A(z), \operatorname{Ad}_{t_{j 1}^{-1}} A(z)\right]\right\rangle \\
& \quad+R_{3}(z, \boldsymbol{\vartheta}) .
\end{align*}
$$

The following is straightforward.
Lemma 4.8.1. The phase (4.90) can be written in the following form

$$
\begin{align*}
\Psi_{x, x}\left(u, \boldsymbol{\vartheta}_{j}, z\right)= & \left(\frac{u}{u_{0}}-1\right)\langle\boldsymbol{\nu}, \boldsymbol{\vartheta}\rangle+u|z|^{2} \sin \left(\beta_{j}-\alpha_{j}\right)\left(\nu_{1}-\nu_{2}\right)  \tag{4.91}\\
& +\frac{i u}{2} \mathbf{w}^{t} G_{x} \mathbf{w}+R_{3}(z, \boldsymbol{\vartheta}),
\end{align*}
$$

where $u_{0}:=\|\Phi(m)\| /\|\boldsymbol{\nu}\|$ and $G_{x}$ is a positive-defined symmetric matrix defined by

$$
\mathbf{w}:=(z, \boldsymbol{\vartheta}) \mapsto\left\|\operatorname{Ad}_{h_{m_{x}}}\left(\left(\operatorname{Ad}_{t_{j 1}^{-1}}-I_{2}\right) A(z)-i \boldsymbol{\vartheta}\right)_{X}(x)\right\|_{x}^{2}
$$

The phase (4.91) has a unique non degenerate critical point,

$$
P_{0}=\left(u_{0}, \mathbf{0}, \mathbf{0}\right),
$$

and the Hessian at the critical point is

$$
H_{P_{0}}\left(\Psi_{x, x}\right)=i\left(\begin{array}{cc}
0 & -i \mathbf{b}_{\nu}^{t} \\
-i \mathbf{b}_{\nu} & H^{\prime}
\end{array}\right)
$$

where $\mathbf{b}_{\boldsymbol{\nu}}^{t}=\left(0,0 ; \boldsymbol{\nu} / u_{0}\right)$ and

$$
i H^{\prime}:=i u_{0} G_{x}+2 u_{0} \sin \left(\beta_{j}-\alpha_{j}\right)\left(\nu_{1}-\nu_{2}\right)\left(\begin{array}{cc}
I_{2} & \mathbf{0} \\
\mathbf{0}^{t} & 0
\end{array}\right) .
$$

We are now in position to apply the Stationary Phase Lemma in (4.88). By Lemma 4.8.1 we have
$\operatorname{det}\left(\frac{H_{P_{0}}\left(\Psi_{j}\right)}{2 \pi i}\right)=-\frac{1}{(2 \pi)^{5}} \operatorname{det}\left(\left(\begin{array}{cc}0 & \mathbf{b}_{\nu}^{t} \\ \mathbf{b}_{\nu} & H^{\prime}\end{array}\right)\right)=\frac{1}{(2 \pi)^{5}} \mathbf{b}_{\nu}^{t}\left(H^{\prime}\right)^{-1} \mathbf{b}_{\nu} \cdot \operatorname{det}\left(H^{\prime}\right)$,
where, for the last equality, we have used Lemma 2.1.3 of [Pao12]. For ease of notation let us set $\Theta_{j}(\boldsymbol{\nu}, x):=\mathbf{b}_{\boldsymbol{\nu}}^{t}\left(H^{\prime}\right)^{-1} \mathbf{b}_{\boldsymbol{\nu}} \cdot \operatorname{det}\left(H^{\prime}\right)$.

Since $H^{\prime}$ and $\left(H^{\prime}\right)^{-1}$ are complex matrices with positive defined real part, we have

$$
\Re\left(\mathbf{b}_{\nu}^{t}\left(H^{\prime}\right)^{-1} \mathbf{b}_{\nu}\right)>0 \quad \text { and } \quad \Re\left(\operatorname{det}\left(H^{\prime}\right)\right)>0
$$

Thus, for both factors appearing in (4.92) we can choose a well-defined square root with argument in $(-\pi / 4, \pi / 4)$, in such a way that

$$
\Re\left(\sqrt{\operatorname{det}\left(\frac{H_{P_{0}}\left(\Psi_{j}\right)}{2 \pi i}\right)}\right)>0 .
$$

Thus, we can compute the leading term of $\Pi_{k \nu}(x, x)_{j 1}, g_{j} \in G_{x} \backslash Z_{x}$ by the Stationary Phase Lemma, it is

$$
\begin{equation*}
\sqrt{2}\left(\frac{\|\boldsymbol{\nu}\|}{\left\|\Phi_{G}\left(m_{x}\right)\right\|}\right)^{d} \cdot D_{G / T} \cdot\left(\frac{k}{\pi}\right)^{d-1 / 2} \cdot \frac{\left(\nu_{1}-\nu_{2}\right)\left(e^{i \alpha}-e^{i \beta}\right)}{\sqrt{\Theta_{j}(\boldsymbol{\nu}, x)}} \cdot e^{-i k\left(\nu_{1} \alpha+\nu_{2} \beta\right)} . \tag{4.93}
\end{equation*}
$$

Notice that $\Pi_{k \nu}(x, x)_{j 2}$ has the same expression as (4.88) with $h_{m_{x}}$ replaced by $k_{m_{x}}$ and $t_{j_{1}}$ replaced by $t_{j_{2}}$. Thus it is an oscillatory integral with phase

$$
\begin{aligned}
& \Psi_{x, x}^{\prime}\left(u, \boldsymbol{\vartheta}_{j}, g T\right)=\frac{i u}{2}\left\|\operatorname{Ad}_{k_{m_{x}}}(\gamma(z, \boldsymbol{\vartheta}))_{X}(x)\right\|_{x}^{2}-\langle\boldsymbol{\nu}, \boldsymbol{\vartheta}\rangle \\
& \quad+u\left\langle\operatorname{Ad}_{k(w)^{-1}}(\Phi(m)), i \boldsymbol{\vartheta}+\frac{1}{2}\left[A(z), \operatorname{Ad}_{t_{j 2}^{-1}} A(z)\right]\right\rangle+R_{3}(z, \boldsymbol{\vartheta}) .
\end{aligned}
$$

In a similar way as before, one can prove that $\Psi_{x, x}^{\prime}$ has not critical point, thus $\Pi_{k \nu}(x, x)_{j 2}$ has rapidly decreasing asymptotic expansion.

Eventually, let us remark that $g_{j} \neq g_{j}^{-1}$ since $g_{j} \notin Z_{x}$; summing the contribution corresponding to $g_{j^{\prime}}:=g_{j}^{-1}$, we obtain the desired asymptotic expansion.

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