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# Transfer operators and horocycle averages on closed manifolds 

Alexander Adam

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## École Doctorale de Sciences Mathématiques de Paris Centre

## Thèse DE DOCTORAT

Discipline: Mathématiques
présentée par

## Alexander Adam

## Opérateurs de transfert et moyennes horocycliques sur les variétés fermées

dirigée par Viviane Baladi

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"To live in an asymmetric world, you should better be symmetric. However if you live in a symmetric world, asymmetry suffices."

L'author

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# Opérateurs de transfert et moyennes horocycliques sur les variétés fermées 

## Résumé

Cette thèse de doctorat approfondit l'étude de la dynamique hyperbolique sur les variétés fermées et connexes $M$ et des opérateurs de transfert associés. Nous étudions deux problèmes : le premier problème concerne les perturbations analytiques réelles des difféomorphismes d'Anosov linéaires sur le tore : une résonance non triviale apparaît-t-elle pour une perturbation génériques d'un difféomorphisme d'Anosov linéaire sur le tore ?
Le second problème concerne une hypothèse sur la moyenne temporelle des flots horocycliques induits par un flot d'Anosov : la moyenne temporelle des flots horocycliques en courbure négative variable converge-t-elle vers la moyenne ergodique en vitesse polynomiale?
Les opérateurs de transfert associés agissent de façon bornée sur certains espaces de Banach anisotropes par la composition du système dynamique inverse suivie d'une multiplication avec des fonctions de poids spécifiques. Dans notre analyse des problèmes mentionnés ci-dessus, ces opérateurs de transfert représentent le principal intérêt. Nous devons étudier leur spectre bas pour progresser sur nos deux problèmes. Par le spectre bas, nous entendons la partie du spectre qui se situe entre le spectre périphérique et le spectre essentiel de ces opérateurs de transfert.

L'approche fonctionnelle de ces opérateurs de transfert se concentre sur les espaces de Banach anisotropes. Nous expliquons l'idée principale derrière cette approche dans le cas des difféomorphismes d'Anosov : des exemples simples de difféomorphismes d'Anosov $F$ sont donnés par les difféomorphismes linéaires d'Anosov sur le tore bidimensionnel. Nous savons que les difféomorphismes d'Anosov transitifs et analytiques ont une unique mesure $\operatorname{SRB} \mu_{\mathrm{SRB}}$ (qui est invariante par le difféomorphisme). Pour les automorphismes linéaires sur le tore, la mesure SRB est la mesure de Lebesgue $\mu_{\text {Leb }}$. Notons toutefois que même de petites perturbations analytiques de $A$ ne préservent pas systématiquement $\mu_{\text {Leb }}$. Puisque $\mu_{\text {SRB }}$ est une mesure de Borel, on a $\mu_{\mathrm{SRB}} \in C(M)^{\prime}$. Nous souhaitons maintenant écrire $\mu_{\mathrm{SRB}}$ comme l'unique vecteur propre associé à la valeur propre 1 pour un certain opérateur de transfert $\mathcal{L}$ qui apparait comme l'adjoint
de l'opérateur de composition $\mathcal{K}_{F}$. Cependant les mesures supportées sur les orbites périodiques de $F$ sont également contenues dans $C(M)^{\prime}$. Afin de trouver les bonnes propriétés spectrales de l'opérateur $\mathcal{L}$, celui-ci doit être défini sur un espace de Banach anisotrope $\mathcal{B}$ et non sur $C(M)^{\prime}$. La norme de $\mathcal{B}$ prend en compte le comportement dilatant et contractant de l'application $F$. En particulier, la norme anisotrope de $\mathcal{B}$ traite les éléments de $\mathcal{B}$ comme des fonctions dans les directions dilatantes et comme des distributions dans les directions contractantes de $F$.

Les valeurs propres discrètes réciproques de $\mathcal{L}$ sont aussi appelées les résonances de $F$. Si $F=A$, alors il y a seulement les résonances triviales $\{0,1\}$. Jusque là il n'était pas su qu'il s'agissait d'un comportement attendu si $A$ est perturbé de manière générique.

On entend ici par perturbation générique toute application d'un ensemble ouvert et dense dans une boule de difféomorphismes analytiques réels contenant $A$. Dans de l'étude du premier problème, nous agissons avec $\mathcal{L}$ sur un espace de Hilbert anisotrope. Nous répondons à la question dans le premier problème par l'affirmative.

Le second problème que nous examinons fait intervenir les flots d'Anosov. Ces flots ont été instaurés par Anosov pour étudier le flot géodésique sur le fibré tangent unitaire de variétés fermées à courbure sectionnelle négative variable. De plus, nous avons besoin les flots d'Anosov d'être des flots de contact. Des exemples de flots d'Anosov-contact sont donnés par les flots géodésiques. Les flots horocycliques associés au flot d'Anosov sont dirigés dans la direction contractant du flot d'Anosov. Nous savons par les travaux de Marcus que pour tout flot horocyclique continu qui correspond à un flot d'Anosov $C^{2}$ mélangeant, il existe une unique mesure de probabilité de Borel invariante par le flot horocyclique. Katok et Burns ont démontré que tout flot d'Anosov-contact est mélangeant. Par conséquent, dans notre contexte, la moyenne temporelle de l'horocycle converge vers la moyenne ergodique unique. Mais à quelle vitesse converge la moyenne temporelle?
Dans le contexte de la courbure négative constante, on sait grâce aux travaux de Flaminio et Forni que cette vitesse est polynomiale. La vitesse est contrôlée par des valeurs propres pour certaines distributions propres du flot géodésique. Un problème analogue dans lequel le flot géodésique est remplacé par un difféomorphisme d'Anosov a été étudié plus tard par Giulietti et Liverani. Ils ont, de plus, dans leurs travaux, supposé que le résultat de Flaminio-Forni devrait s'étendre au flot géodésique dans le contexte de la courbure négative variable.

Dans l'étude du second problème des opérateurs de transfert pondérés $\mathcal{L}_{\alpha}, \alpha>0$ apparaissent. Suivant l'approche fonctionnelle, il suffit essentiellement de construire un espace de Banach anisotrope $\mathcal{B}$ tel que les opérateurs $\mathcal{L}_{\alpha}$ agissent sur $\mathcal{B}$, et d'avoir un spectre périphérique consistant en une valeur propre simple isolée. Cependant, la direction d'écoulement du flot d'Anosov n'est ni contractée ni dilatée par le flot d'Anosov, ce qui pose problème dans notre analyse. Nous appliquons donc à la place la stratégie suivante:
Sur un espace de Banach anisotrope $\mathcal{B}$ bien choisi, la famille d'opérateur de transfert $\left\{\mathcal{L}_{\alpha}: \mathcal{B} \rightarrow \mathcal{B} \mid \alpha \geqslant 0\right\}$ forme un semi-groupe fortement continu et admet donc un générateur bien défini. La quasi-compacité de la résolvante de ce générateur garantit qu'une partie du spectre du générateur par rapport à $\mathcal{B}$ ne contient qu'un spectre discret de multiplicité finie. On trouve alors que le spectre discret contrôle la vitesse de convergence. Cependant, pour montrer que la vitesse de convergence est polynomiale, nous devons imposer une condition supplémentaire à la norme de la résolvante.

Cette thèse est organisée comme suit : dans le premier chapitre, nous étudions les perturbations analytiques réelles et génériques $\widetilde{A}$ d'un difféomorphisme d'Anosov linéaire sur le tore bidimensionnel. Nous appliquons l'approche fonctionnelle aux opérateurs de transfert comme décrit ci-dessus dans le cadre d'un espace de Hilbert anisotrope.
Cet espace de Hilbert est explicitement construit comme la complétion d'un espace de Hardy par rapport à une norme anisotrope. Nous montrons que l'opérateur de composition $\mathcal{K}_{\widetilde{A}}$ est à trace. En fait, nous montrons que $\mathcal{K}_{\tilde{A}}$ est nucléaire d'ordre 0 , ce qui est un résultat plus fort. On obtient la réponse au premier problème par un calcul de la trace de $\mathcal{K}_{\tilde{A}}$ associé au système perturbé. Nous abordons ensuite des propriétés spectrales de l'opérateur de transfert $\mathcal{L}$. Le premier chapitre figure tel qu'il a été publié en 2017 dans Nonlinearity 30.3 , à l'exception de l'annexe qui est séparée en Appendix $A$ et de la liste de références qui est combinée avec celle de cette thèse.

Dans le second chapitre, nous traitons des flots d'Anosov différentiables finis sur des variétés fermées et connexes ainsi que les flots horocycliques stables associés.

La définition de cônes et d'hyperbolicité via des cônes d'une application est développée. Nous introduisons des familles d'opérateurs de transfert $\left\{\mathcal{L}_{\alpha}, \alpha>0\right\}$ avec des fonctions de poids positives arbitraires. Ceci est suivi par la construction d'espaces de Banach anisotropes. Ces espaces sont analogues aux
espaces construits par Baladi et Tsujii pour étudier les difféomorphismes hyperboliques.
Nous étudions les propriétés des opérateurs $\mathcal{L}_{\alpha}, \alpha>0$, sur les espaces de Banach anisotropes ainsi que les propriétés de la résolvante des générateurs des familles $\left\{\mathcal{L}_{\alpha}, \alpha>0\right\}$. Nous établissons une inégalité de Lasota-Yorke pour la résolvante. Nous introduisons et abordons la condition supplémentaire sur les bornes de la résolvante.
Ensuite, nous nous focalisons sur les flots d'Anosov de contact en dimension 3 (nous considérons aussi la co-dimension 1). Nous donnons des bornes locales à l'intégrale de l'horocycle, ce qui nous permet de obtenir notre décomposition de l'intégrale de l'horocycle. Enfin, la condition supplémentaire sur la résolvante est utilisée pour obtenir la vitesse polynomiale de convergence vers la moyenne ergodique.
Le troisième chapitre est également disponible sur arXiv.

## Mots-clés

Flot d'Anosov, flot horocyclique, opérateur de transfert, espace de Banach anisotrope, resonances, moyenne ergodique.

# Transfer operators and horocycle averages on closed manifolds 


#### Abstract

This doctoral thesis deepens the study of hyperbolic dynamics on connected, closed Riemannian manifolds $M$ and associated transfer operators. We investigate two problems: The first problem concerns real analytic perturbations of linear toral Anosov diffeomorphisms: Does a non-trivial resonance appear for generic perturbations of a linear toral Anosov diffeomorphism? The second problem is to make a statement about the time average of horocycle flows with underlying contact Anosov flow: Does the time average of horocycle flows in variable negative curvature converge to the ergodic mean in polynomial time? The associated transfer operators act boundedly on certain anisotropic Banach spaces by composition of the inverse dynamical system followed by a multiplication with specific weight functions. In our analysis of the beforementioned problems these transfer operators are of central interest. We need to investigate their deeper spectrum to progress on our two problems. By the deeper spectrum we mean here the part of the spectrum which lies in between the peripheral and the essential spectrum of these transfer operators. The functional approach to these transfer operators puts importance on the anisotropic Banach spaces. We explain the principal idea behind this approach in the case of Anosov diffeomorphisms: Simple examples of Anosov diffeomorphisms $F$ are provided by the linear Anosov diffeomorphisms $A$ on the twodimensional torus. Real analytic transitive Anosov diffeomorphisms are known to have associated a so-called unique SRB-measure $\mu_{\text {SRB }}$ which is invariant by the diffeomorphism. For the linear toral automorphisms the SRB-measure is just the Lebesgue measure $\mu_{\text {Leb }}$. Note however that even small real-analytic perturbations of $A$ may not preserve $\mu_{\text {Leb }}$. Since $\mu_{\mathrm{SRB}}$ is a Borel measure it holds $\mu_{\mathrm{SRB}} \in C(M)^{\prime}$. We wish now to recover $\mu_{\mathrm{SRB}}$ as the unique 1-eigenvector for a certain transfer operator $\mathcal{L}$ which arises as the adjoint of the composition operator $\mathcal{K}_{F}$. However measures supported on periodic orbits of $F$ are also contained in $C(M)^{\prime}$. In order to find good spectral properties of the operator $\mathcal{L}$, it has to be defined on an anisotropic Banach space $\mathcal{B}$ and not on $C(M)^{\prime}$. The norm of


$\mathcal{B}$ takes into account the expansive and contractive behavior of the map $F$. In particular, the anisotropic norm of $\mathcal{B}$ treats elements in $\mathcal{B}$ as functions along the expanding directions and as distributions along the contracting directions of $F$. The reciprocal discrete eigenvalues of $\mathcal{L}$ are also called the resonances of $F$. In case that $F=A$ there are only the trivial resonances $\{0,1\}$. It was not known before whether this is an expected behavior if $A$ is perturbed generically.

By a generic perturbation we mean here any map in an open and dense subset of a ball of real analytic toral diffeomorphisms containing $A$ with respect to the uniform norm.
In the investigation of the first problem, we act with $\mathcal{L}$ on an anisotropic Hilbert space. We answer the question in the first problem in the affirmative.

The second problem that we investigate involves Anosov flows. These flows where introduced by Anosov to study the geodesic flow on the unit tangent bundle of closed Riemannian manifolds with variable negative sectional curvature. Moreover, we require the contact assumption for the Anosov flow. Examples of contact Anosov flows are provided by geodesic flows on the unit tangent space. The horocycle flows associated to the contact Anosov flow point into the contracting direction of the Anosov flow. We know by the work of Marcus that for every continuous horocycle flow which corresponds to a $C^{2}$ mixing Anosov flow, there exists exactly one Borel probability measure which is invariant by the horocycle flow. Katok and Burns have shown that every contact Anosov flow is mixing. Hence in our setting, the continuous horocycle time average converges to the unique ergodic mean. But how fast is this convergence?
In the constant negative curvature setting, for the geodesic flow, it is known due to the work of Flaminio and Forni, that this speed is polynomial and is controlled by eigenvalues for certain eigendistributions for the geodesic flow. An analogous problem where the geodesic flow is replaced by an Anosov diffeomorphism was studied later by Giulietti and Liverani. Moreover, in their work they conjectured that the above result of Flaminio-Forni should extend to the geodesic flow in variable negative curvature.
In the study of the second problem weighted transfer operators $\mathcal{L}_{\alpha}, \alpha>0$, appear. Following the functional approach, in principal, it is enough to construct an anisotropic Banach space $\mathcal{B}$ such that the operators $\mathcal{L}_{\alpha}$, acting on $\mathcal{B}$, have a peripheral spectrum consisting of an isolated simple eigenvalue. However, the flow direction of the Anosov flow is neither contracted nor expanded by the Anosov flow which poses a problem in our analysis. We apply instead the following strategy:

With a good choice of an anisotropic Banach space $\mathcal{B}$ the transfer operator family $\left\{\mathcal{L}_{\alpha}: \mathcal{B} \rightarrow \mathcal{B} \mid \alpha \geqslant 0\right\}$ forms a strongly continuous semigroup and admits therefore a well-defined generator. Quasi-compactness of the resolvent of this generator ensures that part of the spectrum of the generator with respect to $\mathcal{B}$ contains only discrete spectrum of finite multiplicity. The discrete spectrum is then found to control the speed of convergence. However to show polynomial speed of convergence we need to impose an additional condition on bounds of the resolvent.

## Keywords

Anosov flow, horocycle flow, transfer operator, anistropic Banach space, resonances, ergodic mean.

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## 0 Introduction

This doctoral thesis deepens the study of hyperbolic dynamics on connected, closed Riemannian manifolds $M$ and associated transfer operators.
Two problems are studied: The first problem concerns real analytic perturbations of linear toral Anosov diffeomorphisms. The second problem is to make a detailed statement about the time average of horocycle flows with underlying $C^{3}$ contact Anosov flow. The precise problems are formulated in Problem 1 and Problem 2 in the next section below.
The associated transfer operators act boundedly on certain anisotropic Banach spaces by composition of the inverse dynamical system followed by a multiplication with a weight function. In our analysis of the beforementioned problems for specific weight functions these transfer operators are central objects. We need to investigate their deeper spectral properties to progress on our two problems. By the deeper spectrum we mean here part of the spectrum which lies in between the peripheral and the essential spectrum of these transfer operators.

The functional approach to these transfer operators puts importance on the anisotropic Banach spaces. It is a more recent method with notably results in the last decades for Anosov diffeomorphisms (e.g. [5], [7], [11], [15], [30], [32], [33], [61]) and flows (e.g. [8], [21]-[24], [27], [31], [48], [64], [65]).
We explain the principal idea behind the functional approach briefly in the case of Anosov diffeomorphisms: A diffeomorphic dynamical system $F: M \rightarrow$ $M$ being an Anosov diffeomorphism means that the tangent space $T M$ of the manifold is split into unstable $E_{+} \subset T M$ and stable $E_{-} \subset T M$ distributions. In particular, one assumes a non-trivial splitting

$$
\begin{equation*}
T M=E_{-} \oplus E_{+}, \tag{0.1}
\end{equation*}
$$

such that for every norm $\|\cdot\|$ on linear maps $T M \rightarrow T M$, for some $0<\beta<1$,

$$
\begin{equation*}
\sup _{n \in \mathbb{N}}\left\|\beta^{-n} \mathrm{D} F_{\mid E_{-}}^{n}\right\|<\infty \quad \text { and } \quad \sup _{n \in \mathbb{N}}\left\|\beta^{-n} \mathrm{D} F_{\mid E_{+}}^{-n}\right\|<\infty . \tag{0.2}
\end{equation*}
$$

The arguably simplest example is a toral Anosov diffeomorphism given by Arnold's cat map:

$$
A: \mathbb{R}^{2} / \mathbb{Z}^{2} \rightarrow \mathbb{R}^{2} / \mathbb{Z}^{2}: x \mapsto\left(\begin{array}{ll}
2 & 1 \\
1 & 1
\end{array}\right) x
$$

Similar examples on the torus $\mathbb{R}^{2} / \mathbb{Z}^{2}$ are provided by any hyperbolic unimodular matrix $A \in \mathrm{SL}_{2}(\mathbb{Z})$, where hyperbolic matrix means here that $A$ has one eigenvalue larger than 1 in modulus. Those are the toral linear Anosov diffeomorphisms. The map $A$ is also analytic and mixing. By topological mixing of the map $F$ we mean that for all non-empty subsets $U, V \subseteq M$ there exists $N \in \mathbb{N}$ such that for all $n \geqslant N$ it holds

$$
U \cap F^{n}(V) \neq \varnothing
$$

By transitivity of the map $F$ we mean that there exists a dense orbit of $F$ in $M$. Real analytic transitive Anosov diffeomorphisms $F$ are known to have associated so-called $S R B$-measures $\mu_{\mathrm{SRB}}$ [67, Theorem 1]. (A detailed explanation of $\mu_{\mathrm{SRB}}$ is given in [6], [67].) What is of importance here is that for a given such map $F$ the Borel probability measure $\mu_{\mathrm{SRB}}$ is uniquely characterized by the property that for Lebesgue-almost every $x \in M$ and every continuous function $\varphi \in C(M)$

$$
\mu_{\mathrm{SRB}}(\varphi)=\lim _{n \rightarrow \infty} \frac{1}{n} \sum_{k=0}^{n} \varphi \circ F^{k}(x) .
$$

An example for which $\mu_{\mathrm{SRB}}$ equals the Lebesgue measure $\mu_{\text {Leb }}$ is given by the map $A$. However we note that even small perturbations of $A$ may not preserve $\mu_{\text {Leb }}$.
The composition operator is defined by $\mathcal{K}_{F} \varphi=\varphi \circ F$ for every $\varphi \in C(M)$. Invariance of $\mu_{\text {SRB }}$ by $F$ yields for every $\varphi \in C(M)$

$$
\mu_{\mathrm{SRB}}\left(\mathcal{K}_{F} \varphi\right)=\mu_{\mathrm{SRB}}(\varphi) .
$$

Hence the dual operator $\mathcal{L}:=\mathcal{K}_{F}^{\prime}$ acts on $C(M)^{\prime}$ and fixes $\mu_{\text {SRB }}$. It is straight forward to show that $\mathcal{L}$ with respect to $\mu_{\mathrm{Leb}}$ is given for all $\varphi \in C(M)$ by

$$
\mathcal{L} \varphi=\left|\operatorname{det} \mathrm{D} F^{-1}\right| \cdot \varphi \circ F^{-1}
$$

where $\left|\operatorname{det} \mathrm{D} F^{-1}\right|$ is called the weight function for $\mathcal{L}$. In the functional approach one wishes now to recover $\mu_{\text {SRB }}$ as the unique eigenvector to the eigenvalue 1 for $\mathcal{L}$. More precisely, it is the peripheral spectrum of $\mathcal{L}$ which is here of immediate interest. The peripheral spectrum of $\mathcal{L}$ and the properties of the map $F$ should
be related in the following way:

- The peripheral spectrum of $\mathcal{L}$ contains 1 which is a simple eigenvalue.
- The peripheral spectrum of $\mathcal{L}$ is $\{1\}$ and 1 is a simple eigenvalue.
- The map $F$ is transitive.
- The map $F$ is mixing.

However measures supported on periodic orbits are also contained in $C(M)^{\prime}$ and are therefore eigenvectors to the eigenvalue 1 for $\mathcal{L}$. In order to proceed one changes the domain of the transfer operator $\mathcal{L}$. In particular $\mathcal{L}$ has to be defined on an anisotropic Banach space $\mathcal{B}$ and not on $C(M)^{\prime}$. The space $\mathcal{B}$ contains usually densely $C^{r-1}(M)$ functions for some $r>1$ or $r \in\{\infty, \omega\}$ if $F \in C^{r 1}$. The norm of $\mathcal{B}$ takes into account the expansive and contractive behavior of the map $F$. In particular, the anisotropic norm of $\mathcal{B}$ treats elements in $\mathcal{B}$ as functions along $E_{+}$and as distributions along $E_{-}$.
As we have explained, the construction of $\mathcal{B}$ is constrained by the resulting properties of the spectrum of $\mathcal{L}$ on $\mathcal{B}$. This makes such anisotropic Banach spaces an important part in the functional approach. In the last two decades several constructions of $\mathcal{B}$ have been provided in the differentiable and analytic category of the map $F$ :
The detailed study of anisotropic Banach spaces in the hyperbolic case started in the differentiable setting with the work of Blank, Keller and Liverani [15]. It is now a well established tool (e.g. see the references mentioned above).
Real analytic perturbations of hyperbolic toral automorphisms $A$ were later addressed by Faure and Roy [25]. They considered an anisotropic Hilbert space $\mathcal{H}$, which appeared already briefly in a work of Fried $[29$, Sect $8, I]$.

On this Hilbert space $\mathcal{H}$, the operator $\mathcal{L}$ turns out to be compact which implies that the essential spectrum of $\mathcal{L}$ is contained in $\{0\}$ and the deeper spectrum of $\mathcal{L}$ consist of eigenvalues of finite multiplicity. Moreover the spectrum of $\mathcal{L}$ on $\mathcal{H}$ admits a spectral gap between the deeper and the peripheral spectrum which relates to the following:

- The peripheral spectrum of $\mathcal{L}$ is $\{1\}$ and 1 is an isolated simple eigenvalue.
- The correlation function for $F$ admits exponential decay.

[^0]
## 0 Introduction

By exponential decay of the correlation function for $F$ we mean that for some $\varphi_{1}, \varphi_{2} \in C(M, \mathbb{C})$ and for some $\lambda>0$ it holds

$$
\sup _{n \in \mathbb{N}} e^{\lambda n}\left|\mu_{\mathrm{SRB}}\left(\varphi_{1} \cdot \varphi_{2} \circ F^{n}\right)-\mu_{\mathrm{SRB}}\left(\varphi_{1}\right) \mu_{\mathrm{SRB}}\left(\varphi_{2}\right)\right|<\infty
$$

Usually $\lambda$ depends on the regularity of the observables $\varphi_{1}$ and $\varphi_{2}$ and not on the observables itself.
We comment shortly on the weight function $\left|\operatorname{det} D F^{-1}\right|$ in $\mathcal{L}$ : In general the weight function can be any positive $C^{r-1}$ function, depending on the application. Then of course the maximal eigenvalue $\lambda_{\max }$ may change, as well as the associated eigenvector (this is also called the Gibbs state), replacing $\mu_{\text {SRB }}$. If one considers then the renormalized transfer operator $\widetilde{\mathcal{L}}:=\lambda_{\text {max }}^{-1} \mathcal{L}$, we expect a peripheral spectrum for $\widetilde{\mathcal{L}}$ as discussed, depending on the properties of the map $F$.

The dynamical determinant $d_{F}$ carries statistical information about the behavior of the map $F$ at periodic orbits of $F$. It is a holomorphic function defined for all small enough $z \in \mathbb{C}$ by

$$
d_{F}(z):=\exp -\sum_{n=1}^{\infty} \frac{z^{n}}{n} \sum_{F^{n}(x)=x}\left|\operatorname{det}\left(\mathrm{id}-\mathrm{D}_{x} F^{n}\right)\right|^{-1}
$$

A way to extend the domain of holomorphy of $d_{F}$ is by using the transfer operator. For example if $\mathcal{K}_{F}$ is of trace class then the holomorphic extension is given by the Fredholm determinant for all $z \in \mathbb{C}$ by

$$
\begin{equation*}
\operatorname{det}(1-z \mathcal{L})=\operatorname{det}\left(1-z \mathcal{K}_{F}\right)=d_{F}(z) \tag{0.3}
\end{equation*}
$$

The second equality in (0.3) is a consequence of the trace formula for all $n \in \mathbb{N}$

$$
\begin{equation*}
\operatorname{tr} \mathcal{K}_{F}^{n}=\sum_{F^{n}(x)=x}\left|\operatorname{det}\left(\mathrm{id}-\mathrm{D}_{x} F^{n}\right)\right|^{-1} \tag{0.4}
\end{equation*}
$$

The reciprocal discrete eigenvalues of the transfer operator are also called the (Ruelle-Pollicott) resonances for the map $F$ and the above equality in (0.3) shows a direct relation between the resonances and the zeros of $d_{F}$.

In the hyperbolic setting, Rugh proved the holomorphy of the dynamical determinant of real analytic Anosov diffeomorphisms on surfaces [55], [56].

## The two problems

A small computation shows the following (e.g. using Lemma A.1):
Lemma 1 (Trivial resonances). For every hyperbolic matrix $M \in S L_{2}(\mathbb{Z})$ the associated dynamical determinant satisfies for every $z \in \mathbb{C}$

$$
d_{M}(z)=1-z .
$$

Another direct computation shows:
Lemma 2 (Superexponential decay ([6, Chapter 4.4.1])). For every hyperbolic matrix in $S L_{2}(\mathbb{Z})$ the dynamical correlation function decays superexponentially on real analytic observables.

It was not known before whether the above two lemmas show an expected behavior for a generic real analytic toral Anosov diffeomorphisms.
In particular, what happens if we perturb generically a hyperbolic matrix $A$ by a real analytic map?
By a generic perturbation we mean here any map in an open and dense subset of a ball of real analytic toral diffeomorphisms containing $A$ with respect to the uniform norm (see directly above Theorem 1.4.3).

Problem 1 (Non-trivial resonances). Does a non-trivial resonance appear for generic perturbations of a hyperbolic matrix in $S L_{2}(\mathbb{Z})$ ?

If we consider for a moment only an expanding system $\left(E_{-} \equiv M \times\{0\}\right)$ then more was already known: The expanding case was initially studied by Ruelle [54]. More recently, Bandtlow, Slipantschuk and Just [13], [59] calculated the resonances of real analytic expanding maps $T: S \rightarrow S$ on the unit circle $S$ explicitly for Blaschke products. Their transfer operator acts on the Hardy space of holomorphic functions on the annulus. (See also Keller and Rugh [45] in the differentiable category.)
Moreover, Bandtlow and Naud [12] showed that generically expanding circle maps admit infinitely many resonances.

The second problem that we investigate involves Anosov flows $g_{\alpha} \in C^{3}(M, M)$ where $\operatorname{dim} M=3, \alpha \in \mathbb{R}$ and

$$
g_{\alpha}=g_{\alpha_{1}} \circ g_{\alpha_{2}}, \quad \text { for all } \quad \alpha_{1}+\alpha_{2}=\alpha .
$$

These flows where introduced by Anosov to study the geodesic flow on the unit tangent bundle of closed Riemannian manifolds with variable negative sectional curvature [3], [4]. It is required that the splitting of $T M$ contains in addition a neutral (or central) distribution spanned by the bounded vector field $X$ generating the flow:

$$
\begin{equation*}
T M=E_{+} \oplus E_{-} \oplus \mathbb{R} X \tag{0.5}
\end{equation*}
$$

with an analogous condition on the distributions $E_{-}$and $E_{+}$as in (0.2). Moreover, we require the contact assumption for the flow $g_{\alpha}$ which means that a certain invariant $3-$ form in $\left(T^{*} M\right)^{3}$ is never zero (see Section 2.2). Geodesic flows are well-studied examples of contact Anosov flows.

If $E_{-}$is orientable we define another flow $h_{\rho}: M \rightarrow M, \rho \in \mathbb{R}$, which points into $E_{-}$. This is the (stable) horocycle flow associated to the contact Anosov flow $g_{\alpha}$. (The term horocycle flow was used originally only in the case of the geodesic flow, e.g. see [49, p.84] or [38].) In general the defining vector field of the horocycle flow is at best $C^{2-\epsilon}$ for all $\epsilon>0$ [40].

Theorem 1 (Unique ergodicity, [50, Theorem 3.5]). For every continuous horocycle flow which corresponds to a $C^{2}$ mixing Anosov flow there exists exactly one Borel probability measure which is invariant by the horocycle flow.

Theorem 2 (Mixing ([41, Theorem 3.6], [48, Corollary B.6])). Every contact Anosov flow is mixing.

Hence in our setting, the continuous time average converges to the unique ergodic mean for all $\varphi \in C(M)$ and for all $x \in M$

$$
\mu(\varphi)=\lim _{T \rightarrow \infty} \frac{1}{T} \int_{0}^{T} \varphi \circ h_{\rho}(x) \mathrm{d} \rho
$$

where $\mu$ denotes the unique Borel measure given by Theorem 1. But how fast is the convergence to $\mu(\varphi)$ ? Put otherwise, what can we say about

$$
\frac{1}{T} \int_{0}^{T} \varphi \circ h_{\rho}(x) \mathrm{d} \rho-\mu(\varphi)
$$

for all $T>0$ for fixed $x$ and fixed $\varphi$ ? Clearly, if $\varphi=\partial_{\rho} \widetilde{\varphi} \circ h_{\rho \mid \rho=0}$ for some $\tilde{\varphi} \in C(M, \mathbb{C})$ then

$$
\int_{0}^{T} \varphi \circ h_{\rho}(x) \mathrm{d} \rho=\tilde{\varphi} \circ h_{T}(x)-\tilde{\varphi}(x)
$$

This poses a generic lower bound on the speed of convergence as $\sim 1 / T$, even if $\tilde{\varphi}$ is very regular such that $\varphi \in C^{2-\epsilon}$ for all $\epsilon>0$.
Are there any other obstructions which can slow down the speed of convergence even more? The approach to answer this question is again by means of a weighted transfer operator. We find for every $T, \alpha \geqslant 0$ and every $x \in M$

$$
\gamma_{x}(\varphi, T):=\int_{0}^{T} \varphi \circ h_{\rho}(x) \mathrm{d} \rho=\int_{0}^{\tau(T, \alpha, x)} \mathcal{L}_{\alpha} \varphi \circ h_{\rho} \circ g_{\alpha}(x) \mathrm{d} \rho,
$$

where for all $\alpha \geqslant 0$, assuming here for simplicity that the flow $h_{\rho}$ has unit speed,

$$
\begin{equation*}
\mathcal{L}_{\alpha} \varphi=\operatorname{det} \mathrm{D} g_{-\alpha \mid E_{-}} \cdot \varphi \circ g_{-\alpha}, \tag{0.6}
\end{equation*}
$$

and where for all $x \in M$ and all $\rho, \alpha \in \mathbb{R}$

$$
g_{\alpha} \circ h_{\rho}(x)=h_{\tau(\rho, \alpha, x)} \circ g_{\alpha}(x) .
$$

The function $\tau$ is called the renormalization time. To find a stronger slowdown than $1 / T$ in the speed of convergence, the idea is the following: If on a certain anisotropic Banach space $\mathcal{B}$ the operator $\mathcal{L}_{\alpha}$ has an eigenvector $\mathcal{D}_{\lambda} \in \mathcal{B}$ such that $\mathcal{L}_{\alpha} \mathcal{D}_{\lambda}=e^{\lambda \alpha} \mathcal{D}_{\lambda}$ for some $\Re \lambda>0$ and for all $\alpha \geqslant 0$ then formally

$$
\gamma_{x}\left(\mathcal{D}_{\lambda}, T\right)=\int_{0}^{\tau(T, \alpha, x)} \mathcal{L}_{\alpha} \mathcal{D}_{\lambda} \circ h_{\rho} \circ g_{\alpha}(x) \mathrm{d} \rho=e^{\alpha \lambda} \int_{0}^{\tau(T, \alpha, x)} \mathcal{D}_{\lambda} \circ h_{\rho} \circ g_{\alpha}(x) \mathrm{d} \rho .
$$

As pointed out by Anosov [3], the topological entropy $h_{\text {top }}$ of the time-one map $g_{1}$ is positive. It is an important property of the renormalization time $\tau$ that $\tau(T, \alpha, x) \leqslant 1$ implies $e^{h_{\text {top } \alpha}} \leqslant C T$ for some constant $C>0$ independent of every $T \geqslant 1$ and every $x \in M$ (e.g. use Proposition 2.5.13 below).
Hence in our example, the unique ergodic mean is reached only with at most a speed of $\sim T^{\frac{\Re \lambda}{h_{\text {top }}}-1}$.
Let us assume for simplicity that all eigenvalues $e^{\lambda \alpha}$ are simple for all $\alpha>0$. Then we have associated to each $\lambda$ a finite rank operator given by $\mathcal{D}_{\lambda} \otimes \mathcal{O}_{\lambda}$, where $\mathcal{O}_{\lambda} \in \mathcal{B}^{\prime}$. Now we can decompose formally for every $\varphi \in \mathcal{B}$ and for every $\delta \in \mathbb{R}$

$$
\begin{equation*}
\varphi=\sum_{\Re \lambda>\delta} \mathcal{O}_{\lambda}(\varphi) \mathcal{D}_{\lambda}+\varphi_{\mathcal{E}}, \tag{0.7}
\end{equation*}
$$

with some remainder term $\varphi_{\mathcal{E}}$.

## 0 Introduction

In the constant negative curvature setting, for the geodesic flow, the following is known due to the work of Flaminio and Forni, where the flows $g_{\alpha}$ and $h_{\rho}$ arise both from constant vector fields:

Theorem 3 (Flaminio-Forni, [28, Theorem 1.5]). Let $M$ be the unit tangent bundle of a compact hyperbolic Riemannian surface of constant negative curvature. Let vol be the canonical volume form on $M$. Then it holds for all $\varphi \in C^{4}$, for all $x \in M$ and for all $T>1$

$$
\int_{0}^{T} \varphi \circ h_{\rho}(x) \mathrm{d} \rho=T \operatorname{vol}(\varphi)+\sum_{\lambda \in \Sigma_{\frac{1}{2}} \backslash\{1\}} T^{\lambda} c_{\lambda}(T, x) \mathcal{O}_{\lambda}(\varphi)+\mathcal{E}_{T, x}(\varphi),
$$

where $\sup _{T>1, x \in M} \frac{\left|\mathcal{E}_{T, x}(\varphi)\right|}{T^{\frac{1}{2}} \log T}<\infty$ and $\sup _{T>1, x \in M}\left|c_{\lambda}(T, x)\right|<\infty$.
In fact the result of Flaminio-Forni gives a much more detailed expansion, including a summation over all $\Re \lambda>0$ and lower bounds on the coefficients $c_{\lambda}$.

An analogous problem where the geodesic flow is replaced by an Anosov diffeomorphism was studied later by Giulietti and Liverani [31]. Moreover, in their work they conjectured (see [31, Conjecture 2.14]) that the above result of Flaminio-Forni should extend to the geodesic flow in the variable negative curvature setting.

Problem 2 (Horocycle flows in variable negative curvature). Does an expansion of the horocycle integral analogous to Theorem 3 hold for the horocycle flow induced by the geodesic flow of a surface of variable negative curvature?

Following the functional approach, it is in principle enough to construct an anisotropic Banach space $\mathcal{B}$ such that the operator $\mathcal{L}_{\alpha}$, acting on $\mathcal{B}$, has a peripheral spectrum consisting of an isolated simple eigenvalue at $e^{h_{\text {top }} \alpha}$ for all $\alpha \geqslant 0$ and such that $\gamma_{x}\left(\mathcal{D}_{h_{\text {top }}}, T\right)$ and $\gamma_{x}\left(\varphi_{\mathcal{E}}, T\right)$ are well-defined. Indeed all this could follow if $\mathcal{L}_{\alpha, \phi_{\alpha}}$ is quasi-compact on $\mathcal{B}$ for all $\alpha \geqslant 0$. What prevents us in doing so is the flow direction $X$ in the splitting in (0.5) which is neither contracted nor expanded by the geodesic flow. We apply instead the following strategy: With a good choice of an anisotropic Banach space $\mathcal{B}$ the transfer operator family

$$
\left\{\mathcal{L}_{\alpha}: \mathcal{B} \rightarrow \mathcal{B} \mid \alpha \geqslant 0\right\}
$$

forms a strongly continuous semigroup and admits therefore a well-defined generator $X+V$. Quasi-compactness of the resolvent $\mathcal{R}_{z}$ of $X+V$ for large values
$\Re z>0$ ensures that part of the spectrum of $X+V$ with respect to $\mathcal{B}$

$$
\Sigma_{\delta}:=\left.\sigma(X+V)\right|_{\mathcal{B}} \cap\{\lambda \in \mathbb{C} \mid \Re \lambda>\delta\}
$$

contains only discrete spectrum of finite multiplicity for some $\delta>0$.

Note that $\mathcal{L}_{\alpha}$ is a transfer operator with a different weight than we discussed above in the diffeomorphism case.

In full analogy to our discussion of the peripheral spectrum of $\mathcal{L}$, we expect that the peripheral spectrum of the operator $e^{-h_{\operatorname{top} \alpha}} \mathcal{L}_{\alpha}$ on $\mathcal{B}$ consists of the simple eigenvalue 1 for all $\alpha>0$ and that the associated eigenvector coincides with the unique Borel probability measure given by Theorem 1.

Moreover, the weight function for $\mathcal{L}_{\alpha}$ is $\operatorname{det} \mathrm{D} g_{-\alpha \mid E_{-}}$which depends on the regularity of $E_{-}$, which we noted is $C^{2-\epsilon}$ for all $\epsilon>0$. To deal with such irregular weights one can lift the dynamics to the Grassmanian. This has been used with success, e.g. in [30], [33] and more recently in [65]. However, we handle directly the given weight function which allows us to avoid such additional technicalities.

We should add here that the additional conjecture that the distributions $\mathcal{O}_{v}$ appearing in the right-hand side in (0.7) are fixed by the adjoint of the horocycle flow remains still open. In contrast this was the starting point in [28]. Here, progress has been made by Faure and Guillarmou [24] in dimension 3 for smooth contact Anosov flows.

Although we do not study here the dynamical zeta function for the transfer operator $\mathcal{L}_{\alpha}$, we believe that the anisotropic Banach space $\mathcal{B}$ constructed in Section 2.3.2 could be a suitable choice to be dealt with.

## Statement of results

We present here in a simplified form the main results of this thesis. The first result states that it is quite common for the composition operator to have nontrivial spectrum if the dynamical system is a real analytically perturbed hyperbolic matrix:

Theorem I (Non-trivial resonances (Theorem 1.4.3)). Let $A \in \mathrm{SL}_{2}(\mathbb{Z})$ be hyperbolic. For a generic real analytic perturbation $\tilde{A}$ of $A$ there exists an anisotropic

Hilbert space $\mathcal{H}$ such that

$$
\mathcal{K}_{\tilde{A}}: \mathcal{H} \rightarrow \mathcal{H}
$$

is of trace class and it holds

$$
\operatorname{sp}\left(\mathcal{K}_{\tilde{A}}\right) \backslash\{0,1\} \neq \varnothing
$$

This result implies that the dynamical determinant $d_{\tilde{A}}$ has at least one nontrivial zero and therefore answers the question in Problem 1 in the affirmative. By construction the Hilbert space $\mathcal{H}$ contains real analytic observables as a dense subset. Hence the existence of a non-trivial resonance poses an obstruction to the superexponential decay of the correlation function for real analytic observables.

In some sense our result is optimal: We cannot hope to replace 'generic perturbations' with 'for all perturbations'. In the expanding case there exists Blaschke products arbitrary close to a linear function on the circle and which have trivial spectrum (and which are note $C^{1}$ conjugated to a linear function) [59, Example 5.6].

It is reasonable to believe that similar constructions work in the Anosov case, using the generalized Blaschke products in [61]. This means that isospectral perturbations are expected but they are not generic.

The second result makes a statement about the time average of a horocycle flow with underlying contact Anosov flow. The full result treats the general case with possible non-trivial Jordan blocks. For simplicity of the statement we assume here that the spectrum is simple:

Proposition II (Horocycle integral (Theorem 2.5.7, Proposition 2.5.10)). There exists an anisotropic Banach space $\mathcal{B}$ and $\lambda_{\min }<h_{\text {top }}$ such that $\Sigma_{\lambda_{\text {min }}}$ consists only of discrete eigenvalues of finite multiplicity. For all $x \in M$ and all $T \geqslant 1$, for every $\lambda_{\min } \leqslant \delta<h_{\text {top }}$ and for every finite subset $\Lambda_{\delta} \subseteq \Sigma_{\delta}$ and for all $\varphi \in C^{3}$ it holds
$\int_{0}^{T} \varphi \circ h_{\rho}(x) \mathrm{d} \rho=c_{h_{\text {top }}}(T, x) \mu(\varphi)+\sum_{\substack{\lambda \in \Lambda_{\delta} \\ \Re \lambda<h_{\text {top }}}} T^{\frac{\lambda}{h_{\text {top }}}} c_{\lambda}(T, x) \mathcal{O}_{\lambda}(\varphi)+\mathcal{E}_{T, x, \Lambda_{\delta}}(\varphi)$, where $\mu$ is the unique Borel measure which is invariant by the horocycle flow $h_{\rho}$ and

$$
\sup _{T \geqslant 1, x \in M} \frac{c_{\lambda}(T, x)}{T}<\infty \quad \text { and } \quad \lim _{T \rightarrow \infty} \frac{\mathcal{E}_{T, x, \Lambda_{\delta}}(\varphi)}{T}=0 .
$$

Note that the expected principal term $T \mu(\varphi)$ is obscured since we ordered the expansion by $\mathcal{O}_{\lambda}(\varphi)$. One recovers the principal term, using

$$
T \mu(\varphi)=\int_{0}^{T} 1 \circ h_{\rho}(x) \mathrm{d} \rho .
$$

We use this later in Corollary 2.5.9 in Section 2.5 below. Clearly, we only answer partially the question in Problem 2 in the affirmative. The reason is that quasi-compactness of $\mathcal{R}_{z}$ is not enough to give us e.g. a finite set $\Sigma_{\delta}$ for some $\delta \geqslant \lambda_{\text {min }}$. We impose the following extra condition on the resolvent to affirm the full question for $C^{3}$ contact Anosov flows in dimension 3:

Condition A (Spectral gap with (Dolgopyat) bounds (Condition 2.4.11)). For some $0<\delta<h_{\text {top }}, a>0, b>0, C>0$ and some

$$
\gamma \in\left(0,1 / \log \left(1+\left(h_{\text {top }}-\delta\right) / a\right)\right),
$$

and for all $z \in \mathbb{C}$ with $\Re z=a$ and $|\Im z| \geqslant b$, it holds

$$
\left\|\mathcal{R}_{z+\lambda_{\text {top }}}^{\widetilde{n}}\right\|_{W_{p}^{s, t, q}} \leqslant C^{\widetilde{n}}\left|\Re z+\left(h_{\text {top }}-\delta\right)\right|^{-\widetilde{n}}, \quad \text { where } \tilde{n}=\lceil\gamma \log |\Im z|\rceil
$$

Under this additional condition we obtain:

Theorem III (Theorem 2.5.7, Proposition 2.5.10). Under the assumptions of Proposition II, if in addition Condition A holds with same $\delta$ then we can take $\Lambda_{\delta}=\Sigma_{\delta}$ and it holds $\left|\mathcal{E}_{T, x, \Sigma_{\delta}}(\varphi)\right| \leqslant C T^{\frac{\delta}{h_{\text {top }}}+\epsilon}$ for all $\epsilon>0$ and some $C=$ $C(\varphi) \geqslant 0$ independent of $T$ and $x$.

We shall note a curiosity which we do not discuss further in this thesis: An application of the last theorem which presents itself is the deeper analysis of the renormalization time $\tau$ itself! It follows from the construction of the transfer operator $\mathcal{L}_{\alpha}$ that for all $\rho, \alpha>0$ and for all $x \in M$ it holds

$$
\tau(\rho,-\alpha, x)=\gamma_{x}\left(\mathcal{L}_{\alpha} 1, \rho\right) .
$$

Of course, in the setting of constant vector fields (i.e. constant negative curvature) we cannot learn anything new about $\tau$, but other cases might be of further interest (e.g. in studying small perturbations of the constant vector fields).

## Organization of the thesis

We investigate Problem 1 in Chapter 1 and Problem 2 in Chapter 2.

In Chapter 1 we study generic real analytic perturbations $\tilde{A}$ of a linear Anosov diffeomorphism on the two-dimensional torus as introduced before Problem 1. We apply the functional approach to transfer operators as described above in the setting of an anisotropic Hilbert space.
This Hilbert space is explicitly constructed in Section 1.2 as the completion of some Hardy space with respect to an anisotropic norm. The trace class property of the Koopman operator $\mathcal{K}_{\widetilde{A}}$ is shown in Section 1.3. In fact, we show that $\mathcal{K}_{\widetilde{A}}$ is nuclear of order 0 which is a stronger result.

Theorem I which is Theorem 1.4.3 below is finally shown in Section 1.4. This is done essentially by the calculation of the trace of the transfer operator associated to the perturbed system, using the equality in (0.4) which is shown as well. The spectral properties of the transfer operator $\mathcal{L}$ are discussed in Section 1.5. Chapter 1 is presented as it was published [1], except that the appendix is moved to Chapter A and the reference list is combined with that of this thesis.

In Chapter 2 we deal with finite differentiable Anosov flows on connected, closed Riemannian manifolds and the associated stable horocycle flows.
The necessary notion of cones and cone-hyperbolicity of a map is introduced in Section 2.2.

In Section ?? we introduce families $\left\{\mathcal{L}_{\alpha} \mid \alpha>0\right\}$ of transfer operators similar to those in (0.6) but with arbitrary positive weight functions. This is followed by the construction of anisotropic Banach spaces. These spaces are a flow analogue to the spaces constructed by Baladi and Tsujii [10] to study hyperbolic diffeomorphisms.

In Section 2.4 we discuss properties of the operators $\mathcal{L}_{\alpha}, \alpha>0$, on the constructed anisotropic Banach spaces as well as properties of the resolvent of the generators of the families $\left\{\mathcal{L}_{\alpha} \mid \alpha>0\right\}$. We show a Lasota-Yorke inequality for the resolvent, which is Theorem 2.4.5 below. We introduce and discuss Condition A which is Condition 2.4.11 below.

In Section 2.5 we specialize to contact Anosov flows in dimension 3 (we consider also the co-dimension 1 case).
We give local bounds on the horocycle integral in Lemma 2.5.14. This lemma enables us to show Theorem III which is Theorem 2.5 .7 below.

The additional Condition A is finally needed to obtain a polynomial rate of
convergence to the ergodic mean for horocycle flows induced by contact Anosov flows in dimension 3.

This chapter is also available on arXiv [2].

## 1 Generic non-trivial resonances for Anosov diffeomorphisms

### 1.1 Introduction

Let $T: \mathbb{T}^{2} \rightarrow \mathbb{T}^{2}$ be a real analytic Anosov diffeomorphism. We define the Ruelle resonances of $T$ to be the zeroes of the (holomorphically continued in $z \in \mathbb{C}$ ) dynamical determinant

$$
\begin{equation*}
d_{T}(z):=\exp -\sum_{n=1}^{\infty} \frac{z^{n}}{n} \sum_{T^{n}(x)=x}\left|\operatorname{det}\left(\mathrm{id}-\mathrm{D}_{x} T^{n}\right)\right|^{-1} . \tag{1.1}
\end{equation*}
$$

It is well-known (e.g. combining (1.1) and Lemma A.1) that 1 is the only resonance if $T$ is a hyperbolic linear toral automorphism $M$. A subset of the Banach space of $\mathbb{T}^{2}$-preserving maps, holomorphic and uniformly bounded on some annulus, is called generic if it is open and dense. We show in Theorem 1.4.3, using an idea of Naud [51], that there is such a set $\mathcal{G}$ so that for all $\psi \in \mathcal{G}$, appropriately scaled, the Anosov diffeomorphism $M+\psi$ admits non-trivial Ruelle resonances. For this, we construct a Hilbert space of anisotropic generalized functions on which the transfer operator $\mathcal{L}_{T} f:=(f /|\operatorname{det} \mathrm{D} T|) \circ T^{-1}$ is nuclear with its Fredholm determinant equal to $d_{T}$. Moreover, we prove that some of those generic perturbations preserve the volume while some do not.
The expanding case is easier and was initially studied by Ruelle [54]. More recently, Bandtlow et. al [13], [59] calculated the resonances of real analytic expanding maps $T: S \rightarrow S$ on the unit circle $S$ explicitly for Blaschke products. Their transfer operator acts on the Hardy space of holomorphic functions on the annulus. (See also Keller and Rugh [45] in the differentiable category.)
In the hyperbolic setting, Rugh proved the holomorphy of the dynamical determinant of real analytic Anosov diffeomorphisms on surfaces [55], [56]. The idea was generalized by Fried to hyperbolic flows in all dimensions [29]. The detailed study of anisotropic Banach spaces in the hyperbolic case started with the pioneering work of [15] (in the differentiable setting) and is now a well established
tool, see e.g. [11] and [32].
Faure and Roy [25] later addressed real analytic perturbations of hyperbolic linear toral automorphisms on the two-dimensional torus, considering an anisotropic complex Hilbert space, which had already been briefly discussed by Fried [29, Sect 8, I].

Our approach is based on this construction and strongly relies on an idea suggested by Naud [51]. We put the transfer operator at the center of our analysis. We introduce an anisotropic Hilbert space (Definition 1.2.4) in Section 1.2.
In Section 1.3, we rephrase a result from Faure and Roy [25, Theorem 6] to show that the Koopman operator $\mathcal{K}_{T} f:=f \circ T$ is nuclear of order 0 when acting on our anisotropic Hilbert space.
In Section 1.4, we use this result and an idea of Naud [51] to show that the Koopman operator admits non-trivial Ruelle resonances under a small generic perturbation of the dynamics.
In Section 1.5, we consider the adjoint of the Koopman operator, which is just the transfer operator, acting on the dual Hilbert space and obtain our final results.
In the Appendix, we recall two needed basic properties of integer matrices (seen as linear maps on the torus) and provide a sufficient condition for determinant preserving perturbations of differentiable real maps.

In principal the analogous problem on any higher dimensional torus can be treated with the presented method. However, one has to modify slightly the used space from Section 1.2 if the linear toral automorphism has non-trivial Jordan blocks.

Blaschke products were recently generalized to the hyperbolic setting by Slipantschuk et al. [60] who calculate the entire spectrum of these real analytic Anosov volume preserving diffeomorphisms explicitly.

### 1.2 An anisotropic Hilbert space

We denote the flat 2 -torus by $\mathbb{T}^{2}:=\mathbb{R}^{2} / \mathbb{Z}^{2}$. We embed $\mathbb{T}^{2}$ into the standard polyannulus in $\mathbb{C}^{2}$ and set for each $r>0$

$$
\mathbb{A}_{r}:=\mathbb{T}^{2}+\mathrm{i}(-r, r)^{2}
$$

We see $\mathbb{A}_{r}$ as a submanifold of $\mathbb{C}^{2}$. The Hilbert space $L_{2}\left(\mathbb{T}^{2}\right)$ is equipped with the canonical Lebesgue measure on $\mathbb{T}^{2}$. This space admits an orthonormal Fourier basis given by

$$
\begin{equation*}
\varphi_{n}: \mathbb{T}^{2} \rightarrow \mathbb{C}: x \mapsto \exp \left(\mathrm{i} 2 \pi n^{*} x\right), \quad n \in \mathbb{Z}^{2} \tag{1.2}
\end{equation*}
$$

where $n^{*}$ is the canonical dual of $n$. We recall a construction from Faure and Roy [25] for a complex Hilbert space $\mathcal{H}_{A_{M, c}}$. This space also has been described briefly by Fried as an "ad hoc example" [29, Sect. 8, I.] of a generalized function space. The construction will be based on:
Definition 1.2.1 (Hardy space $\left.H_{2}\left(\mathbb{A}_{r}\right)\right)$. For each $r>0$ and each holomorphic function $f: \mathbb{A}_{r} \rightarrow \mathbb{C}$, we define the norm

$$
\|f\|_{H_{2}\left(\mathbb{A}_{r}\right)}:=\sup _{y \in(-r, r)^{2}}\left(\int_{\mathbb{T}^{2}}|f(x+\mathrm{i} y)|^{2} \mathrm{~d} x\right)^{\frac{1}{2}} .
$$

Then we set

$$
H_{2}\left(\mathbb{A}_{r}\right):=\left\{f: \mathbb{A}_{r} \rightarrow \mathbb{C} \mid f \text { holomorphic, }\|f\|_{H_{2}\left(\mathbb{A}_{r}\right)}<\infty\right\}
$$

The space $H_{2}\left(\mathbb{A}_{r}\right)$ is the 2-dimensional analogue of the Hardy space studied in [58, p. 4]. It admits a Fourier basis given by

$$
\vartheta_{n}^{r}: \mathbb{A}_{r} \rightarrow \mathbb{C}: x \mapsto \exp (-2 \pi r\|n\|) \varphi_{n}, \quad n \in \mathbb{Z}^{2}
$$

where $\|z\|:=\left|z_{1}\right|+\left|z_{2}\right|$ for all $\left(z_{1}, z_{2}\right)=: z \in \mathbb{C}^{2}$ and $z \in \mathbb{T}^{2}$. With this choice of norm, the Fourier basis is orthonormal. Under the canonical isomorphism $L_{2}\left(\mathbb{T}^{2}\right) \cong L_{2}\left(\mathbb{T}^{2}\right)^{*}$, we have the isomorphism

$$
\begin{equation*}
\left(\vartheta_{n}^{r}\right)^{*} \cong \vartheta_{n}^{-r} \tag{1.3}
\end{equation*}
$$

A matrix $M \in \mathrm{SL}_{2}(\mathbb{Z})$ is called hyperbolic if its eigenvalues do not lie on the unit circle. We denote by $E_{M}^{+}$the eigenspace for the eigenvalue of modulus $\lambda_{M}>1$ and by $E_{M}^{-}$the eigenspace of the eigenvalue of modulus $\lambda_{M}^{-1}$. We decompose $y \in \mathbb{R}^{2}$ uniquely as

$$
\begin{equation*}
y=y_{M}^{+}+y_{M}^{-} \quad \text { with } \quad y_{M}^{+} \in E_{M^{*}}^{+} \quad \text { and } \quad y_{M}^{-} \in E_{M^{*}}^{-} \tag{1.4}
\end{equation*}
$$

We have

$$
\begin{equation*}
\left\|M^{*} y_{M}^{+}\right\|=\lambda_{M}\left\|y_{M}^{+}\right\| \quad \text { and } \quad\left\|M^{*} y_{M}^{-}\right\|=\lambda_{M}^{-1}\left\|y_{M}^{-}\right\| \tag{1.5}
\end{equation*}
$$

1 Generic non-trivial resonances for Anosov diffeomorphisms

Definition 1.2.2 (Scaling map $\left.A_{M, c}\right)$. Let $c>0$, and $M \in \mathrm{SL}_{2}(\mathbb{Z})$ be hyperbolic. For every $n \in \mathbb{Z}^{2}$, we set, recalling (1.2),

$$
A_{M, c} \varphi_{n}:=\exp \left(-2 \pi c\left(\left\|n_{M}^{+}\right\|-\left\|n_{M}^{-}\right\|\right)\right) \varphi_{n} .
$$

Lemma 1.2.3 (Continuous embedding of $H_{2}\left(\mathbb{A}_{r}\right)$ ). Let $c>0$ and let $M \in$ $\mathrm{SL}_{2}(\mathbb{Z})$ be hyperbolic. Then the map $A_{M, c}$ can be extended by continuity to an injective linear map

$$
A_{M, c}: H_{2}\left(\mathbb{A}_{c}\right) \rightarrow L_{2}\left(\mathbb{T}^{2}\right),
$$

bounded in operator norm by 1.
Proof. By Definition 1.2.2, for each $f \in H_{2}\left(\mathbb{A}_{c}\right)$ we have

$$
\begin{aligned}
\left\|A_{M, c} f\right\|_{L_{2}\left(\mathbb{T}^{2}\right)}^{2} & =\sum_{n \in \mathbb{Z}^{2}}\left|\varphi_{n}^{*} A_{M, c} f\right|^{2}=\sum_{n \in \mathbb{Z}^{2}} \exp \left(-4 \pi c\left(\left\|n_{M}^{+}\right\|-\left\|n_{M}^{-}\right\|\right)\right)\left|\varphi_{n}^{*} f\right|^{2} \\
& =\sum_{n \in \mathbb{Z}^{2}} \exp \left(-4 \pi c\left(\left\|n_{M}^{+}\right\|-\left\|n_{M}^{-}\right\|+\|n\|\right)\right)\left|\vartheta_{n}^{c *} f\right|^{2},
\end{aligned}
$$

where we used (1.3) in the last step. Using the triangle inequality, we find

$$
\left\|n_{M}^{+}\right\|-\left\|n_{M}^{-}\right\|+\|n\| \geqslant 0 .
$$

Hence, it holds

$$
\sum_{n \in \mathbb{Z}^{2}} \exp \left(-4 \pi c\left(\left\|n_{M}^{+}\right\|-\left\|n_{M}^{-}\right\|+\|n\|\right)\right)\left|\vartheta_{n}^{c *} f\right|^{2} \leqslant\|f\|_{H_{2}\left(\mathbb{A}_{c}\right)}^{2} .
$$

Injectivity follows since $A_{M, c}$ is invertible on the Fourier basis of $L_{2}\left(\mathbb{T}^{2}\right)$.
The image of $H_{2}\left(\mathbb{A}_{c}\right)$ under $A_{M, c}$ is dense in $L_{2}\left(\mathbb{T}^{2}\right)$ since it contains all Fourier polynomials.

Definition 1.2.4 (Hilbert space $\left.\mathcal{H}_{A_{M, c}}\right)$. Let $c>0$ and let $M \in \mathrm{SL}_{2}(\mathbb{Z})$ be hyperbolic. Let $A_{M, c}$ be the map given by Definition 1.2.2. Then we set

$$
\mathcal{H}_{A_{M, c}}:=\text { closure of } H_{2}\left(\mathbb{A}_{c}\right) \text { with respect to the norm }\left\|A_{M, c} \cdot\right\|_{L_{2}\left(\mathbb{T}^{2}\right)} \text {, }
$$

and extend $A_{M, c}$ by continuity to a linear map

$$
A_{M, c}: \mathcal{H}_{A_{M, c}} \rightarrow L_{2}\left(\mathbb{T}^{2}\right) .
$$

As a direct consequence of this construction, the scalar product on $\mathcal{H}_{A_{M, c}}$ satisfies

$$
\langle\cdot, \cdot\rangle_{\mathcal{H}_{A_{M, c}}}: \mathcal{H}_{A_{M, c}} \times \mathcal{H}_{A_{M, c}} \rightarrow \mathbb{C}:(f, g) \mapsto\left\langle A_{M, c} f, A_{M, c} g\right\rangle_{L_{2}\left(\mathbb{T}^{2}\right)}
$$

An orthonormal Fourier basis of $\mathcal{H}_{A_{M, c}}$ is given by

$$
\begin{equation*}
\varrho_{n}:=A_{M, c}^{-1} \varphi_{n}, \quad n \in \mathbb{Z}^{2} . \tag{1.6}
\end{equation*}
$$

Lemma 1.2.5 (Dual space of $\left.\mathcal{H}_{A_{M, c}}\right)$. Under the canonical isomorphism $L_{2}\left(\mathbb{T}^{2}\right) \cong$ $L_{2}\left(\mathbb{T}^{2}\right)^{*}$, the dual space $\mathcal{H}_{A_{M, c}}^{*}$ is isomorphic to $A_{M, c}^{2} \mathcal{H}_{A_{M, c}}$.

Proof. Under the canonical isomorphism $L_{2}\left(\mathbb{T}^{2}\right) \cong L_{2}\left(\mathbb{T}^{2}\right)^{*}$, we have for each $n_{1}, n_{2} \in \mathbb{Z}^{2}$, using (1.6),

$$
\varphi_{n_{1}}^{*}\left(\varphi_{n_{2}}\right)=\varphi_{n_{1}}^{*}\left(A_{M, c} \varrho_{n_{2}}\right)=\left(A_{M, c} \varphi_{n_{1}}\right)^{*}\left(\varrho_{n_{2}}\right)=\left(A_{M, c}^{2} \varrho_{n_{1}}\right)^{*}\left(\varrho_{n_{2}}\right) .
$$

Remark 1.2.6. By Lemma 1.2.5, we associate to every linear functional f* $\in$ $\mathcal{H}_{A_{M, c}}^{*}$ a unique vector $f \in A_{M, c}^{2} \mathcal{H}_{A_{M, c}}$. Then, for every $g \in \mathcal{H}_{A_{M, c}}$, the product $f g$ is absolutely integrable with respect to the Lebesgue measure on $\mathbb{T}^{2}$.

The decomposition in (1.4) defines two cones

$$
C_{M}^{+}:=\left\{y \in \mathbb{R}^{2} \mid\left\|y_{M}^{+}\right\| \geqslant\left\|y_{M}^{-}\right\|\right\} \quad \text { and } \quad C_{M}^{-}:=\left\{y \in \mathbb{R}^{2} \mid\left\|y_{M}^{+}\right\| \leqslant\left\|y_{M}^{-}\right\|\right\}
$$

Example 1.2.7. We let $M=\left(\begin{array}{ll}3 & 1 \\ 2 & 1\end{array}\right)$, then $\lambda_{M}=2+\sqrt{3}$. An eigenvector for $\lambda_{M}$ for $M^{*}$ is $(1+\sqrt{3}, 1)$ and an eigenvector for $\lambda_{M}^{-1}$ is $(1-\sqrt{3}, 1)$. The two subspaces $E_{M^{*}}^{+}$and $E_{M^{*}}^{-}$and the two cones $C_{M}^{+}$and $C_{M}^{-}$are shown in Figure 1.1.

We set

$$
\begin{aligned}
& \mathcal{H}_{A_{M, c}}^{+}:=\left\{\sum_{n \in C_{M}^{+} \cap \mathbb{Z}^{2}}\left\langle\varrho_{n}, f\right\rangle_{\mathcal{H}_{A_{M, c}}} \varrho_{n} \mid f \in \mathcal{H}_{A_{M, c}}\right\} \text { and } \\
& \mathcal{H}_{A_{M, c}}^{-}:=\left\{\sum_{n \in C_{M}^{-} \cap \mathbb{Z}^{2}}\left\langle\varrho_{n}, f\right\rangle_{\mathcal{H}_{A_{M, c}}} \varrho_{n} \mid f \in \mathcal{H}_{A_{M, c}}\right\} .
\end{aligned}
$$

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Figure 1.1: The map $M$ is from Example 1.2.7. The dark gray area is the cone $C_{M}^{+}$which contains the subspace $E_{M^{*}}^{+}$. The light gray area is the cone $C_{M}^{-}$and contains $E_{M^{*}}^{-}$. A part $y \in \mathbb{R}^{2}$ belongs to the dashed lines if and only if $\left\|y_{M}^{+}\right\|=\left\|y_{M}^{-}\right\|$.

Hence, we have $\mathcal{H}_{A_{M, c}}=\mathcal{H}_{A_{M, c}}^{+}+\mathcal{H}_{A_{M, c}}^{-}$. Comparing for each $n \in C_{M}^{-}$the Fourier basis $\varrho_{n}$ with $\varphi_{n}$, it follows immediately that $\mathcal{H}_{A_{M, c}}^{-} \subset L_{2}\left(\mathbb{T}^{2}\right)$. For each $n \in C_{M}^{+}$, comparing the Fourier basis $\varrho_{n}$ with $\vartheta_{n}^{c *}$, using (1.3), shows $\mathcal{H}_{A_{M, c}}^{+} \subset H_{2}\left(\mathbb{A}_{c}\right)^{*}$. We conclude therefore that $\mathcal{H}_{A_{M, c}}$ contains linear functionals which do not belong to $L_{2}\left(\mathbb{T}^{2}\right)$. By construction, the space $\mathcal{H}_{A_{M, c}}$ is a rigged Hilbert space, i.e.:

$$
\begin{equation*}
H_{2}\left(\mathbb{A}_{c}\right) \subset \mathcal{H}_{A_{M, c}} \subset H_{2}\left(\mathbb{A}_{c}\right)^{*} . \tag{1.7}
\end{equation*}
$$

Remark 1.2.8. We note that in the construction of $\mathcal{H}_{A_{M, c}}$, the expanding and contracting directions appear in the dual coordinates $n \in \mathbb{Z}^{2}$ of the Fourier basis (1.6). This distinguishes $\mathcal{H}_{A_{M, c}}$ from the space of Rugh [56] where expanding and contracting coordinates are spatial. We observe

$$
n^{*} x=\left(n_{M}^{+}+n_{M}^{-}\right)^{*}\left(x_{M^{*}}^{+}+n_{M^{*}}^{-}\right)=\left(n_{M}^{+}\right)^{*} x_{M^{*}}^{+}+\left(n_{M}^{-}\right)^{*} x_{M^{*}}^{-} .
$$

Hence, we can rewrite (1.6) as

$$
\begin{align*}
\varrho_{n}(x)= & \exp \left(2 \pi c\left(\left\|n_{M}^{+}\right\|-\left\|n_{M}^{-}\right\|\right)\right) \exp \left(\mathrm{i} 2 \pi n^{*} x\right) \\
= & \exp \left(2 \pi c\left\|n_{M}^{+}\right\|\right) \exp \left(\mathrm{i} 2 \pi\left(n_{M}^{+}\right)^{*} x_{M^{*}}^{+}\right)  \tag{1.8}\\
& \quad \times \exp \left(-2 \pi c\left\|n_{M}^{-}\right\|\right) \exp \left(\mathrm{i} 2 \pi\left(n_{M}^{-}\right)^{*} x_{M^{*}}^{-}\right) .
\end{align*}
$$

It is tempting to think of the $\varrho_{n}$ as basis elements for a tensor product space of
a Hardy space on an annulus, with the dual of such a Hardy space. However, we cannot use $\varrho_{n}$ as such a basis since $n_{M}^{+}$and $n_{M}^{-}$are not independent of each other. Nevertheless, we can decompose $\mathcal{H}_{A_{M, c}}$ into two generalized Hardy spaces as follows. We define four norms

$$
\begin{aligned}
\mu_{j}(f) & :=\sup _{y \in A_{j}}\left(\int_{\mathbb{T}^{2}}|f(x+\mathrm{i} y)|^{2} \mathrm{~d} x\right)^{\frac{1}{2}}, f \in L_{2}\left(\mathbb{T}^{2}\right), j \in\{1,2,3,4\}, \text { where } \\
A_{1} & :=\left\{y \in \mathbb{R}^{2} \mid y_{M^{*}}^{-} \in(-c, c)^{2}, y_{M^{*}}^{+} \in(c, \infty)^{2}\right\}, \\
A_{2} & :=\left\{y \in \mathbb{R}^{2} \mid y_{M^{*}}^{-} \in(-c, c)^{2}, y_{M^{*}}^{+} \in(-\infty,-c)^{2}\right\}, \\
A_{3} & :=\left\{y \in \mathbb{R}^{2} \mid y_{M^{*}}^{-} \in(-c, c)^{2}, y_{M^{*}}^{+} \in(c, \infty) \times(-\infty,-c)\right\}, \\
A_{4} & :=\left\{y \in \mathbb{R}^{2} \mid y_{M^{*}}^{-} \in(-c, c)^{2}, y_{M^{*}}^{+} \in(-\infty,-c) \times(c, \infty)\right\} .
\end{aligned}
$$

For all $f \in L_{2}\left(\mathbb{T}^{2}\right)$ the norms $\mu_{j}(f)$ cannot be finite but they are so at least for some Fourier polynomials. The spaces $H_{j}, j \in\{1,2,3,4\}$, are the completions with respect to the norms $\mu_{j}$ above. E.g. using $\mu_{1}$, it holds for all $f \in H_{1}$

$$
\begin{aligned}
\mu_{1}(f)^{2} & =\sup _{y \in A_{1}}\left(\int_{\mathbb{T}^{2}}|f(x+\mathrm{i} y)|^{2} \mathrm{~d} x\right)=\sup _{y \in A_{1}} \sum_{n \in \mathbb{Z}^{2}} \exp \left(-4 \pi n^{*} y\right)\left|\varphi_{n}^{*} f\right|^{2} \\
& =\sup _{y \in A_{1}} \sum_{n \in \mathbb{Z}^{2}} \exp \left(-4 \pi\left(n_{M}^{-}\right)^{*} y_{M^{*}}^{-}-4 \pi\left(n_{M}^{+}\right)^{*} y_{M^{*}}^{+}\right)\left|\varphi_{n}^{*} f\right|^{2} \\
& =\sup _{y_{M^{*}}^{+} \in(c, \infty)^{2}} \sum_{n \in \mathbb{Z}^{2}} \exp \left(4 \pi c\left\|n_{M}^{-}\right\|-4 \pi\left(n_{M}^{+}\right)^{*} y_{M^{*}}^{+}\right)\left|\varphi_{n}^{*} f\right|^{2} \\
& =\sum_{\substack{n \in \mathbb{Z}^{2} \\
n_{M}^{+} \in[0, \infty)^{2}}} \exp \left(4 \pi c\left\|n_{M}^{-}\right\|-4 \pi c\left\|n_{M}^{+}\right\|\right)\left|\varphi_{n}^{*} f\right|^{2}=\sum_{\substack{n \in \mathbb{Z}^{2} \\
n_{M}^{+} \in[0, \infty)^{2}}}\left|\varphi_{n}^{*} A_{M, c} f\right|^{2} .
\end{aligned}
$$

Similar calculations for the other three norms show then that the spaces $H_{j}, j \in$ $\{1,2,3,4\}$ disjointly partition the space $\mathcal{H}_{A_{M, c}}$ with respect to the dual coordinate up to $n=0$. Since $E_{M}^{+}$is a one dimensional subspace of $\mathbb{R}^{2}$, always two of the spaces contain only the constant functions (note that $n_{M}^{+}=0$ implies $n=0$ ), say, $H_{3}$ and $H_{4}$. Then all vectors in the spaces $H_{1}$ and $H_{2}$ are holomorphic functions on $\mathbb{T}^{2}+\mathrm{i} A_{1}$ and on $\mathbb{T}^{2}+\mathrm{i} A_{2}$, respectively.

### 1.3 The Koopman operator is nuclear

We set for each $r>0$
$\mathcal{T}_{r}:=\left\{T: \mathbb{T}^{2} \rightarrow \mathbb{T}^{2} \mid T\right.$ extends holomorphically and boundedly on $\left.\mathbb{A}_{r}\right\}$.

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For every $T \in \mathcal{T}_{r}$ the Koopman operator

$$
\mathcal{K}_{T}: L_{2}\left(\mathbb{T}^{2}\right) \rightarrow L_{2}\left(\mathbb{T}^{2}\right): f \mapsto f \circ T
$$

is well-defined by differentiability of $T$. It is well-known that the operator $\mathcal{K}_{T}$ acting on $L_{2}\left(\mathbb{T}^{2}\right)$ is not compact. We say that two maps $f, g \in \mathcal{T}_{r}$ are $C^{1}$-close if the distance

$$
d(f, g):=\sup _{z \in \mathbb{A}_{r}}\|f(z)-g(z)\|+\sup _{z \in \mathbb{A}_{r}}\left\|\mathrm{D}_{z} f-\mathrm{D}_{z} g\right\|
$$

is small. In this section we revisit the proof of Faure and Roy [25]. They showed that $\mathcal{K}_{T}$, acting on the Hilbert space $\mathcal{H}_{A_{M, c}}$, (see Definition 1.2.4), is nuclear of order 0 if $T$ is sufficiently $C^{1}$-close to a hyperbolic matrix $M \in \mathrm{SL}_{2}(\mathbb{Z})$ for some $c>0$.
We recall that a linear operator $\mathcal{L}: \mathcal{H} \rightarrow \mathcal{H}$ on a Hilbert space $\mathcal{H}$ with norm $\|\cdot\|_{\mathcal{H}}$ is called nuclear of order 0 if it can be written as a sum $\mathcal{L}=\sum_{n \in \mathbb{N}} d_{n} \psi_{1, n} \psi_{2, n}^{*}$ with $\inf \left\{p>\left.0\left|\sum_{n \in \mathbb{N}}\right| d_{n}\right|^{p}<\infty\right\}=0$ and $\psi_{1, n}, \psi_{2, n} \in \mathcal{H},\left\|\psi_{1, n}\right\|_{\mathcal{H}},\left\|\psi_{2, n}\right\|_{\mathcal{H}} \leqslant 1$, $d_{n} \in \mathbb{C}, n \in \mathbb{N}\left[34, \mathrm{II}, \S 1, \mathrm{n}^{\circ} 1, \mathrm{p} .4\right]$. In particular, such an operator is trace class, hence bounded and admits a $\operatorname{trace} \operatorname{tr} \mathcal{L}:=\sum_{n \in \mathbb{N}} e_{n}^{*} \mathcal{L} e_{n}$, invariant for any choice of orthonormal basis $e_{n}, n \in \mathbb{N}$ of $\mathcal{H}$. Moreover, one can show that $\operatorname{tr} \mathcal{L}$ equals the sum, including multiplicity (dimension of corresponding generalized eigenspace), over the spectrum $\operatorname{sp}(\mathcal{L})$ of $\mathcal{L}$. The Fredholm determinant, defined for small enough $z \in \mathbb{C}$ by

$$
\begin{equation*}
\operatorname{det}(1-z \mathcal{L}):=\exp \left(-\sum_{n=1}^{\infty} \frac{z^{n}}{n} \operatorname{tr} \mathcal{L}^{n}\right) \tag{1.9}
\end{equation*}
$$

extends to an entire function in $z$, having zeroes at $z=\lambda^{-1}, \lambda \in \operatorname{sp}(\mathcal{L}) \backslash\{0\}$ of same order as the multiplicity of $\lambda$.

Theorem 1.3.1 (Nuclearity of $\mathcal{K}_{T}$ ). Let $M \in \mathrm{SL}_{2}(\mathbb{Z})$ be hyperbolic and let $r>0$. Then there exist constants $\delta_{M}>0$ and $0<c_{1}<r$ such that for each $T \in \mathcal{T}_{r}$ with $d(T, M) \leqslant \delta_{M}$ the map

$$
\mathcal{K}_{T}: \mathcal{H}_{A_{M, c_{1}}} \rightarrow \mathcal{H}_{A_{M, c_{1}}}: f \mapsto f \circ T
$$

defines a nuclear operator of order 0 . In particular, there exists $c_{2}>0$ depending
only on $c_{1}, M$, and $\|\cdot\|$ so that for each $n_{1}, n_{2} \in \mathbb{Z}^{2}$

$$
\left|\left\langle\varrho_{n_{1}}, \mathcal{K}_{T} \varrho_{n_{2}}\right\rangle_{\mathcal{H}_{A_{M, c_{1}}}}\right| \leqslant \exp \left(-2 \pi c_{2}\left(\left\|n_{1}\right\|+\left\|n_{2}\right\|\right)\right) .
$$

For every $n_{1}, n_{2} \in \mathbb{Z}^{2}$, we set

$$
\begin{equation*}
I_{n_{1}, n_{2}}(T):=\left\langle\varphi_{n_{1}}, \mathcal{K}_{T} \varphi_{n_{2}}\right\rangle_{L_{2}\left(\mathbb{T}^{2}\right)} \tag{1.10}
\end{equation*}
$$

Estimating this "oscillatory integral" is central for Theorem 1.3.1. In the case $T=M$, we have simply

$$
I_{n_{1}, n_{2}}(M)= \begin{cases}1 & \text { if } M^{*} n_{2}=n_{1}  \tag{1.11}\\ 0 & \text { if } M^{*} n_{2} \neq n_{1}\end{cases}
$$

The strategy of the proof is as follows. We get an upper bound for $\left|I_{n_{1}, n_{2}}(T)\right|$ in Lemma 1.3.2, taking advantage of the holomorphicity of $T$. In Lemma 1.3.3, we compare the contribution of $n_{1}$ and $n_{2}$ in the expanding and contracting directions, using here essentially the hyperbolicity of $M$. Combining both results, we obtain a weaker bound on $\left|I_{n_{1}, n_{2}}(T)\right|$ in Proposition 1.3.4, which finally allows for the proof of Theorem 1.3.1.
For every $n \in \mathbb{Z}^{2}$ and $y \in \mathbb{R}^{2}$ any solution $x \in \mathbb{T}^{2}$ so that

$$
\begin{equation*}
\exp \left(-2 \pi\left(n^{*} \mathrm{D}_{x} T y\right)\right)=\int_{\mathbb{T}^{2}} \exp \left(-2 \pi\left(n^{*} \mathrm{D}_{z} T y\right)\right) \mathrm{d} z \tag{1.12}
\end{equation*}
$$

is denoted by $x_{n}(y)$. Since the integrand is continuous in $y$ such a solution exists by the Mean Value Theorem.

Lemma 1.3.2 (Upper bound on $\left|I_{n_{1}, n_{2}}(T)\right|(\mathrm{I})$ ). Let $r>0$. Then, there exists $C \geqslant 0$ so that for each $T \in \mathcal{T}_{r}$ and $n_{1}, n_{2} \in \mathbb{Z}^{2}$ and $y \in(-r, r)^{2}$, recalling (1.10), we have

$$
\left|I_{n_{1}, n_{2}}(T)\right| \leqslant \exp \left(2 \pi\left(-n_{2}^{*} \mathrm{D}_{x_{n_{2}}(y)} T y+n_{1}^{*} y+C d(T, 0)\|y\|^{3}\left\|n_{2}\right\|\right)\right)
$$

Proof. By definition

$$
I_{n_{1}, n_{2}}(T)=\left\langle\varphi_{n_{1}}, \mathcal{K}_{T} \varphi_{n_{2}}\right\rangle_{L_{2}\left(\mathbb{T}^{2}\right)}=\int_{\mathbb{T}^{2}} \exp \left(\mathrm{i} 2 \pi\left(n_{2}^{*} T(x)-n_{1}^{*} x\right)\right) \mathrm{d} x
$$

Since $T \in \mathcal{T}_{r}$, the $\mathbb{Z}^{2}$-invariance of the integrand follows. By holomorphicity of $T$ on $\mathbb{A}_{r}$, we can change the path of integration to $x \mapsto x+\mathrm{i} y$ for every $y \in(-r, r)^{2}$.

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Therefore for any $y \in(-r, r)^{2}$

$$
\left|I_{n_{1}, n_{2}}(T)\right| \leqslant \int_{\mathbb{T}^{2}} \exp \left(2 \pi\left(n_{1}^{*} y-\Im\left(n_{2}^{*} T(x+\mathrm{i} y)\right)\right)\right) \mathrm{d} x
$$

where $\Im$ is the imaginary part. We expand $T$ (or rather its lift to $\mathbb{R}^{2}$ ) at $x \in \mathbb{T}^{2}$ in a Taylor series to the second order. This yields

$$
T(x+\mathrm{i} y)=T(x)+\mathrm{iD}_{x} T y+P(x+\mathrm{i} y)+R_{2}(x+\mathrm{i} y) .
$$

Here, $P(x+\mathrm{i} y)$ is the second order term of the expansion which is $\mathbb{R}^{2}$-valued, and $R_{2}$ is the remainder of the series expansion. We find therefore

$$
\Im T(x+\mathrm{i} y)=\mathrm{D}_{x} T y+\Im R_{2}(x+\mathrm{i} y) .
$$

Since $T$ is holomorphic we find a constant $C>0$ independent of $T$ such that

$$
\left|n_{2}^{*} R_{2}(x+\mathrm{i} y)\right| \leqslant C d(T, 0)\left\|n_{2}\right\|\|y\|^{3} .
$$

We are left with the evaluation of

$$
\int_{\mathbb{T}^{2}} \exp \left(-2 \pi\left(n^{*} \mathrm{D}_{z} T y\right)\right) \mathrm{d} z .
$$

Using (1.12) yields the result.
The following abbreviation is used in the remaining section. We set for each $y \in \mathbb{R}^{2}$

$$
\begin{equation*}
|y|_{M}:=\left\|y_{M}^{+}\right\|-\left\|y_{M}^{-}\right\| . \tag{1.13}
\end{equation*}
$$

Lemma 1.3.3 (Directional inequality). Let $M \in \mathrm{SL}_{2}(\mathbb{Z})$ be hyperbolic. Let $\epsilon>0$ and $\kappa \geqslant 0$ and let $R: \mathbb{R}^{2} \rightarrow \mathbb{R}_{\geqslant 0}$ be a map such that for all $z \in \mathbb{R}^{2}$ with $\|z\|<\epsilon$ it holds

$$
R(z) \leqslant \kappa\|z\| .
$$

Then there exists $c_{M}>0$ such that if $\kappa<c_{M}$ there exist $0<c_{2}<c_{1}<\epsilon$ such that for all $n_{1}, n_{2} \in \mathbb{Z}^{2}$ there exists $y_{n_{1}, n_{2}} \in \mathbb{R}^{2}$ independent of $R$ with $\left\|y_{n_{1}, n_{2}}\right\|<\epsilon$ such that it holds
$-c_{1}\left(\left|n_{1}\right|_{M}-\left|n_{2}\right|_{M}\right)-\left(n_{2}^{*} M-n_{1}^{*}\right) y_{n_{1}, n_{2}}+\left\|n_{2}\right\| R\left(y_{n_{1}, n_{2}}\right) \leqslant-c_{2}\left(\left\|n_{1}\right\|+\left\|n_{2}\right\|\right)$.
Proof. We assume $0<c_{2} \leqslant c_{1}$. For $n_{1}=n_{2}=0$ there is nothing to prove. For every $\left(y_{1}, y_{2}\right) \in \mathbb{R}^{2}$ we set $\left|\left(y_{1}, y_{2}\right)\right|:=\sqrt{y_{1}^{2}+y_{2}^{2}}$. We let $0<\tilde{c}_{1} \leqslant 1 \leqslant \tilde{c}_{2}$ such
that

$$
\begin{equation*}
\tilde{c}_{2}^{-1}\left(\left\|y_{M}^{+}\right\|+\left\|y_{M}^{-}\right\|\right) \leqslant\|y\| \leqslant \tilde{c}_{1}^{-1}|y|, \quad \text { for all } y \in \mathbb{R}^{2} \tag{1.14}
\end{equation*}
$$

Whenever $n_{2} \neq 0$ we find a linear map $M_{a}$ such that $M_{a} n_{2}=M^{*} n_{2}-n_{1}$ and whenever $n_{1} \neq 0$ we find a linear map $M_{b}$ such that $M_{b} n_{1}=M^{*} n_{2}-n_{1}$. For now we let $\tilde{\kappa}>0$ be a variable which will be fixed later on, independently of $n_{1}$ and $n_{2}$. We consider the following four cases
(a) $\left\|n_{2}\right\|>0$ and $\left\|n_{2}\right\| \geqslant\left\|n_{1}\right\|$
(b) $\left\|n_{1}\right\|>0$ and $\left\|n_{1}\right\| \geqslant\left\|n_{2}\right\|$
(i) $\left\|M_{a} n_{2}\right\| \geqslant \tilde{\kappa}\left\|n_{2}\right\|$,
(i) $\left\|M_{b} n_{1}\right\| \geqslant \tilde{\kappa}\left\|n_{1}\right\|$,
(ii) $\left\|M_{a} n_{2}\right\|<\tilde{\kappa}\left\|n_{2}\right\|$,
(ii) $\left\|M_{b} n_{1}\right\|<\tilde{\kappa}\left\|n_{1}\right\|$.

We assume Case (a)(i). For every $\delta>0$ we let

$$
y=\delta M_{a} \frac{n_{2}}{\left\|n_{2}\right\|}
$$

It follows, using (1.14), that

$$
\begin{equation*}
-\left(n_{2}^{*} M-n_{1}^{*}\right) y=-n_{2}^{*} M_{a}^{*} y \leqslant-\tilde{c}_{1}^{2}\left\|M_{a} n_{2}\right\|\|y\| \tag{1.15}
\end{equation*}
$$

We recall $|\cdot|_{M}$ from (1.13). Using that $c_{1}+c_{2}>0$ and that (a) holds, we estimate

$$
\begin{aligned}
-c_{1}\left(\left|n_{1}\right|_{M}-\left|n_{2}\right|_{M}\right) & \leqslant c_{1}\left(\left\|n_{1}\right\|+\left\|n_{2}\right\|\right) \\
& =-c_{2}\left(\left\|n_{1}\right\|+\left\|n_{2}\right\|\right)+\left(c_{1}+c_{2}\right)\left(\left\|n_{1}\right\|+\left\|n_{2}\right\|\right) \\
& \leqslant-c_{2}\left(\left\|n_{1}\right\|+\left\|n_{2}\right\|\right)+2\left(c_{1}+c_{2}\right)\left\|n_{2}\right\|
\end{aligned}
$$

Using (a)(i) and the assumed bound on $R$ for $\|y\|<\epsilon$, we have

$$
\begin{align*}
-c_{1}\left(\left|n_{1}\right|_{M}-\left|n_{2}\right|_{M}\right) & -\left\|n_{2}\right\|\left(\tilde{c}_{1}^{2}\left\|M_{a} \frac{n_{2}}{\left\|n_{2}\right\|}\right\|\|y\|-R(y)\right) \leqslant \\
& -c_{2}\left(\left\|n_{1}\right\|+\left\|n_{2}\right\|\right)+\left(2\left(c_{1}+c_{2}\right)-\left(\tilde{c}_{1}^{2} \tilde{\kappa}-\kappa\right)\|y\|\right)\left\|n_{2}\right\| \tag{1.16}
\end{align*}
$$

We put $c_{M}:=\tilde{c}_{1}^{2} \tilde{\kappa}$. Any value $\|y\| \in(0, \epsilon)$ can be attained by controlling $\delta$. Assuming that $c_{M}>\kappa$, it follows from (1.15) and (1.16) that

$$
\begin{equation*}
0<c_{1}+c_{2}<\frac{c_{M}-\kappa}{2} \epsilon \tag{1.17}
\end{equation*}
$$

The reasoning in Case (b)(i) is completely analogous and yields the same bounds

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on $c_{1}+c_{2}$.
In Case (a)(ii) and (b)(ii), we take $y=0$, where $R(0)=0$ by assumption on $R$.
We assume now Case (a)(ii). We find, using (1.14),

$$
\begin{align*}
\left\|\left(M_{a} n_{2}\right)_{M}^{+}\right\|+\left\|\left(M_{a} n_{2}\right)_{M}^{-}\right\| & \leqslant \tilde{c}_{2}\left\|M_{a} n_{2}\right\| \\
& <\tilde{c}_{2} \tilde{\kappa}\left\|n_{2}\right\| \leqslant \tilde{c}_{2} \tilde{\kappa}\left(\left\|n_{2, M}^{+}\right\|+\left\|n_{2, M}^{-}\right\|\right) . \tag{1.18}
\end{align*}
$$

We have

$$
\left\|\left(M_{a} n_{2}\right)_{M}^{+}\right\|=\left\|M^{*} n_{2, M}^{+}-n_{1, M}^{+}\right\| \quad \text { and } \quad\left\|\left(M_{a} n_{2}\right)_{M}^{-}\right\|=\left\|M^{*} n_{2, M}^{-}-n_{1, M}^{-}\right\| .
$$

Recalling (1.5), this allows the estimate

$$
\begin{aligned}
\left\|\left(M_{a} n_{2}\right)_{M}^{+}\right\|+\left\|\left(M_{a} n_{2}\right)_{M}^{-}\right\| & \geqslant\left\|M^{*} n_{2, M}^{+}\right\|-\left\|n_{1, M}^{+}\right\|-\left\|M^{*} n_{2, M}^{-}\right\|+\left\|n_{1, M}^{-}\right\| \\
& \geqslant \lambda_{M}\left\|n_{2, M}^{+}\right\|-\lambda_{M}^{-1}\left\|n_{2, M}^{-}\right\|-\left\|n_{1, M}^{+}\right\|+\left\|n_{1, M}^{-}\right\| .
\end{aligned}
$$

Together with (1.18) we find therefore

$$
-\left|n_{1}\right|_{M}=-\left\|n_{1, M}^{+}\right\|+\left\|n_{1, M}^{-}\right\|<-\left(\lambda_{M}-\tilde{\kappa} \tilde{c}_{2}\right)\left\|n_{2, M}^{+}\right\|+\left(\lambda_{M}^{-1}+\tilde{\kappa} \tilde{c}_{2}\right)\left\|n_{2, M}^{-}\right\| .
$$

We set

$$
\kappa_{+}:=\lambda_{M}-\tilde{\kappa} \tilde{c}_{2}-1 \quad \text { and } \quad \kappa_{-}:=1-\lambda_{M}^{-1}-\tilde{\kappa} \tilde{c}_{2} .
$$

We finally estimate

$$
-c_{1}\left(\left|n_{1}\right|_{M}-\left|n_{2}\right|_{M}\right)<-c_{1} \kappa_{+}\left\|n_{2, M}^{+}\right\|-c_{1} \kappa_{-}\left\|n_{2, M}^{-}\right\| .
$$

Note that we have $\kappa_{+}>\kappa_{-}$because $\lambda_{M}>1$. Assuming that $c_{1} \kappa_{-} \geqslant 2 c_{2}$, we find

$$
-c_{1} \kappa_{-}\left\|n_{2, M}^{+}\right\|-c_{1} \kappa_{-}\left\|n_{2, M}^{-}\right\|<-c_{1} \kappa_{-}\left\|n_{2}\right\| \leqslant-2 c_{2}\left\|n_{2}\right\| \leqslant-c_{2}\left(\left\|n_{1}\right\|+\left\|n_{2}\right\|\right) .
$$

In Case (b)(ii) we consider the bounds

$$
\begin{aligned}
\left|n_{2}\right|_{M}+\lambda_{M}\left\|n_{1, M}^{-}\right\|-\lambda_{M}^{-1}\left\|n_{1, M}^{+}\right\| & \leqslant\left\|\left(\left(M^{*}\right)^{-1} M_{b} n_{1}\right)_{M}^{+}\right\|+\left\|\left(\left(M^{*}\right)^{-1} M_{b} n_{1}\right)_{M}^{-}\right\| \\
& \leqslant \tilde{c}_{2}\left\|\left(M^{*}\right)^{-1} M_{b} n_{1}\right\|<\tilde{\kappa} \tilde{c}_{2}\left\|\left(M^{*}\right)^{-1}\right\|\left\|n_{1}\right\| .
\end{aligned}
$$

Therefore $\kappa_{-}$is replaced by $1-\lambda_{M}^{-1}-\left\|\left(M^{*}\right)^{-1}\right\| \tilde{\kappa} \tilde{c}_{2}$ which we require to be
positive. Since $\left\|\left(M^{*}\right)^{-1}\right\|>1$, this yields the stronger conditions

$$
\begin{equation*}
0<\tilde{\kappa}<\frac{1-\lambda_{M}^{-1}}{\left\|\left(M^{*}\right)^{-1}\right\| \tilde{c}_{2}} \quad \text { and } \quad c_{2} \leqslant \frac{1-\lambda_{M}^{-1}-\left\|\left(M^{*}\right)^{-1}\right\| \tilde{\kappa} \tilde{c}_{2}}{2} c_{1} \tag{1.19}
\end{equation*}
$$

Any such choice for $\tilde{\kappa}$ is independent of $n_{1}$ and $n_{2}$ and fixes $c_{M}$. Using (1.19) for an upper bound on $c_{2}$ and (1.17), we find the stronger condition

$$
0<c_{1}<\frac{c_{M}-\kappa}{3-\lambda_{M}^{-1}} \epsilon
$$

Therefore the choices of $c_{1}$ and $c_{2}$ are valid if $\kappa<c_{M}$. They depend only on $\epsilon$, $M$ and $\|\cdot\|$ and not on $n_{1}$ or $n_{2}$.

Proposition 1.3.4 (Upper bound on $\left|I_{n_{1}, n_{2}}(T)\right|$ (II)). Let $M \in \mathrm{SL}_{2}(\mathbb{Z})$ be hyperbolic and let $r>0$. Then there exist constants $0<\delta_{M}$ and $0<c_{2}<c_{1}<r$ such that for each $n_{1}, n_{2} \in \mathbb{Z}^{2}$ and each $T \in \mathcal{T}_{r}$ with $d(T, M) \leqslant \delta_{M}$ it holds that

$$
\exp \left(-2 \pi c_{1}\left(\left|n_{1}\right|_{M}-\left|n_{2}\right|_{M}\right)\right)\left|I_{n_{1}, n_{2}}(T)\right| \leqslant \exp \left(-2 \pi c_{2}\left(\left\|n_{1}\right\|+\left\|n_{2}\right\|\right)\right)
$$

Proof. By Lemma 1.3.2 there is a constant $C>0$ independent of $T$ such that for each $y \in(-r, r)^{2}$ and $n_{1}, n_{2} \in \mathbb{Z}^{2}$ it holds that

$$
\begin{equation*}
\left|I_{n_{1}, n_{2}}(T)\right| \leqslant \exp \left(2 \pi\left(-n_{2}^{*} \mathrm{D}_{x_{n_{2}}(y)} T y+n_{1}^{*} y+C d(T, 0)\|y\|^{3}\left\|n_{2}\right\|\right)\right) . \tag{1.20}
\end{equation*}
$$

We rewrite

$$
n_{2}^{*} \mathrm{D}_{x_{n_{2}}(y)} T y=n_{2}^{*} M y+n_{2}^{*} \mathrm{D}_{x_{n_{2}}(y)}(T-M) y,
$$

and set

$$
R(y):=\left\{\begin{array}{ll}
\frac{n_{2}^{*}}{\left\|n_{2}\right\|} \mathrm{D}_{x_{n_{2}}(y)}(M-T) y+C d(T, 0)\|y\|^{3} & \text { if } n_{2} \neq 0 \\
0 & \text { if } n_{2}=0
\end{array} .\right.
$$

Let $\delta_{M}>0$ and assume that $d(T, M) \leqslant \delta_{M}$. We choose $0<\epsilon \leqslant r$ sufficiently small such that for all $y \in \mathbb{R}^{2}$ with $\|y\|<\epsilon$ there is $\kappa>0$ such that

$$
|R(y)| \leqslant \kappa \delta_{M} .
$$

Since $d(T, 0) \leqslant d(T, M)+d(M, 0) \leqslant \delta_{M}+d(M, 0)$ this choice of $\epsilon$ is independent

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of $T$. Lemma 1.3.3 applied to $M$ and $|R|$ gives $c_{1}, c_{2}$ and $y_{n_{1}, n_{2}} \in \mathbb{R}^{2}$ for which the right-hand side of (1.20) fulfills the desired inequality.

Proof of Theorem 1.3.1. Proposition 1.3.4 yields $0<\delta_{M}$ and $0<c_{2}<c_{1}<r$ such that if $d(T, M) \leqslant \delta_{M}$ it holds

$$
\begin{equation*}
C_{n_{1}} C_{n_{2}}^{-1}\left|I_{n_{1}, n_{2}}(T)\right| \leqslant \exp \left(-2 \pi c_{2}\left(\left\|n_{1}\right\|+\left\|n_{2}\right\|\right)\right), \tag{1.21}
\end{equation*}
$$

where

$$
C_{n}:=\exp \left(-2 \pi c_{1}\left(\left\|n_{M}^{+}\right\|-\left\|n_{M}^{-}\right\|\right)\right), \quad n \in \mathbb{Z}^{2} .
$$

We put $c:=c_{1}$ and $M$ in Definitions 1.2.2 and 1.2.4, giving a linear map $A_{M, c_{1}}$ and a Hilbert space $\mathcal{H}_{A_{M, c_{1}}}$. Recalling (1.6), and assuming that $\mathcal{K}_{T}: \mathcal{H}_{A_{M, c_{1}}} \rightarrow$ $\mathcal{H}_{A_{M, c_{1}}}$ is well-defined, we have

$$
\begin{align*}
\left|\left\langle\varrho_{n_{1}}, \mathcal{K}_{T} \varrho_{n_{2}}\right\rangle_{\mathcal{H}_{A_{M, c_{1}}}}\right| & =\left|\left\langle\varphi_{n_{1}}, A_{M, c_{1}} \mathcal{K}_{T} A_{M, c_{1}}^{-1} \varphi_{n_{2}}\right\rangle_{L_{2}\left(\mathbb{T}^{2}\right)}\right| \\
& =C_{n_{1}} C_{n_{2}}^{-1}\left|I_{n_{1}, n_{2}}(T)\right| . \tag{1.22}
\end{align*}
$$

Using (1.21) to estimate the right-hand side, the bound in Theorem 1.3.1 follows. We next obtain well-definedness and nuclearity of order 0 of $\mathcal{K}_{T}$. Let $f \in \mathcal{H}_{A_{M, c_{1}}}$ and put $g:=A_{M, c_{1}} f$. We have then

$$
\begin{aligned}
\mathcal{K}_{T} f \in \mathcal{H}_{A_{M, c_{1}}} & \Leftrightarrow A_{M, c_{1}} \mathcal{K}_{T} f \in L_{2}\left(\mathbb{T}^{2}\right) \Leftrightarrow \sum_{n \in \mathbb{Z}^{2}}\left|\varphi_{n}^{*} A_{M, c_{1}} \mathcal{K}_{T} f\right|^{2}<\infty \\
& \Leftrightarrow \sum_{n_{1} \in \mathbb{Z}^{2}}\left|\sum_{n_{2} \in \mathbb{Z}^{2}} \varphi_{n_{1}}^{*} A_{M, c_{1}} \mathcal{K}_{T} A_{M, c_{1}}^{-1} \varphi_{n_{2}} \varphi_{n_{2}}^{*} g\right|^{2}<\infty \\
& \Leftrightarrow \sum_{n_{1} \in \mathbb{Z}^{2}}\left|\sum_{n_{2} \in \mathbb{Z}^{2}} C_{n_{1}} C_{n_{2}}^{-1} I_{n_{1}, n_{2}}(T) \varphi_{n_{2}}^{*} g\right|^{2}<\infty .
\end{aligned}
$$

Using (1.21) and the Cauchy-Schwartz inequality, it follows that

$$
\sum_{n_{1} \in \mathbb{Z}^{2}}\left|\sum_{n_{2} \in \mathbb{Z}^{2}} C_{n_{1}} C_{n_{2}}^{-1} I_{n_{1}, n_{2}}(T) \varphi_{n_{2}}^{*} g\right|^{2} \leqslant\left(\sum_{n \in \mathbb{Z}^{2}} e^{-4 \pi c_{2}\|n\|}\right)^{2}\|g\|_{L_{2}\left(\mathbb{T}^{2}\right)}^{2}<\infty .
$$

This gives the well-definedness of $\mathcal{K}_{T}$. Now, using the Cauchy-Schwartz inequa-
lity, we have

$$
\left|\left\langle\varrho_{n}, \mathcal{K}_{T} f\right\rangle_{\mathcal{H}_{A_{M, c_{1}}}}\right|^{2} \leqslant \sum_{m \in \mathbb{Z}^{2}}\left|\left\langle\varrho_{n}, \mathcal{K}_{T} \varrho_{m}\right\rangle_{\mathcal{H}_{A_{M, c_{1}}}}\right|^{2}\|f\|_{\mathcal{H}_{A_{M, c_{1}}}}^{2}
$$

Using (1.22) and (1.21) to bound $\left|\left\langle\varrho_{n}, \mathcal{K}_{T} \varrho_{m}\right\rangle_{\mathcal{H}_{A_{M, c_{1}}}}\right|$, we find a constant $C>0$ such that

$$
\left|\left\langle C \exp \left(2 \pi c_{2}\|n\|\right) \varrho_{n}, \mathcal{K}_{T} f\right\rangle_{\mathcal{H}_{A_{M, c_{1}}}}\right| \leqslant\|f\|_{\mathcal{H}_{A_{M, c_{1}}}} .
$$

This allows the representation of $\mathcal{K}_{T}$ as

$$
\mathcal{K}_{T} f=\sum_{n \in \mathbb{Z}^{2}} C^{-1} \exp \left(-2 \pi c_{2}\|n\|\right)\left\langle C \exp \left(2 \pi c_{2}\|n\|\right) \varrho_{n}, \mathcal{K}_{T} f\right\rangle_{\mathcal{H}_{A_{M, c_{1}}}} \varrho_{n}
$$

from which nuclearity of order 0 follows. Finally, a brief inspection of the proofs for Lemma 1.3.3 and Proposition 1.3.4 gives the statement about the constants.

### 1.4 Non-trivial resonances for the Koopman operator

Given any hyperbolic matrix $M \in \mathrm{SL}_{2}(\mathbb{Z})$, we find by Theorem 1.3.1 constants $0<\delta_{M}$ and $c>0$ such that for each map $T \in \mathcal{T}_{r}$, satisfying $d(T, M) \leqslant \delta_{M}$, the operator $\mathcal{K}_{T}$ acting on the Hilbert space $\mathcal{H}_{A_{M, c}}$ is nuclear of order 0 . Therefore it has a well-defined trace

$$
\begin{equation*}
\operatorname{tr} \mathcal{K}_{T}:=\sum_{n \in \mathbb{Z}^{2}}\left\langle\varrho_{n}, \mathcal{K}_{T} \varrho_{n}\right\rangle_{\mathcal{H}_{A_{M, c}}} \tag{1.23}
\end{equation*}
$$

The map $T$ is an Anosov diffeomorphism (for all small enough $\delta_{M}$ ), by structural stability [35, Theorem 9.5.8]. Then the map $T$ has the same number $N_{M}=$ $|\operatorname{det}(\operatorname{id}-M)|$ of fixed points as the matrix $M$. We recall a well-known result [25, Proposition 9].

Lemma 1.4.1 (Trace formula for $\left.\mathcal{K}_{T}\right)$. Let $M \in \mathrm{SL}_{2}(\mathbb{Z})$ be hyperbolic and let $r>0$. Then there exist constants $\delta_{M}>0$ and $c>0$ such that for each $T \in \mathcal{T}_{r}$ with $d(T, M) \leqslant \delta_{M}$, letting $\mathcal{K}_{T}$ act on $\mathcal{H}_{A_{M, c}}$, it holds

$$
\operatorname{tr} \mathcal{K}_{T}=\sum_{T(x)=x}\left|\operatorname{det}\left(\mathrm{id}-\mathrm{D}_{x} T\right)\right|^{-1}
$$

For the convenience of the reader, we give a proof:

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Proof. Using Theorem 1.3.1 gives constants $c>0$ and $\delta_{M}>0$ and welldefinedness of $\mathcal{K}_{T}$. For small enough $\delta_{M}>0$, by structural stability and Lemma A. 1 (ii), the map id $-T$ can be partitioned into $N_{M}$ surjective submaps. In particular, there are diffeomorphisms $y_{j}: D_{j} \rightarrow \mathbb{T}^{2}, D_{j} \subseteq \mathbb{T}^{2}, 1 \leqslant j \leqslant N_{M}$ such that id $-T=\bigsqcup_{j=1}^{N_{M}} y_{j}$. Then, using (1.6), we have for each $n \in \mathbb{Z}^{2}$

$$
\begin{aligned}
\left\langle\varrho_{n}, \mathcal{K}_{T} \varrho_{n}\right\rangle_{\mathcal{H}_{A_{M, c}}} & =\left\langle\varphi_{n}, A_{M, c} \mathcal{K}_{T} A_{M, c}^{-1} \varphi_{n}\right\rangle_{L_{2}\left(\mathbb{T}^{2}\right)}=\int_{\mathbb{T}^{2}} \exp \left(\mathrm{i} 2 \pi n^{*}(T-\mathrm{id})(x)\right) \mathrm{d} x \\
& =\sum_{j=1}^{N_{M}} \int_{y_{j}^{-1}\left(\mathbb{T}^{2}\right)} \exp \left(\mathrm{i} 2 \pi n^{*} y_{j}(x)\right) \mathrm{d} x \\
& =\sum_{j=1}^{N_{M}} \int_{\mathbb{T}^{2}} \frac{\exp \left(\mathrm{i} 2 \pi n^{*} z\right)}{\operatorname{det}\left(\operatorname{id}-\mathrm{D}_{y_{j}^{-1}(z)} T\right)} \mathrm{d} z .
\end{aligned}
$$

For $N \in \mathbb{N}$ and $z \in \mathbb{T}^{2}$ the following sum

$$
D_{N}(z):=\sum_{\substack{n \in \mathbb{Z}^{2} \\\|z\| \leqslant N}} \exp \left(\mathrm{i} 2 \pi n^{*} z\right)
$$

is the 2-dimensional analogue of the Dirichlet kernel [44, p.13]. Together with (1.23), this yields immediately

$$
\operatorname{tr} \mathcal{K}_{T}=\lim _{N \rightarrow \infty} \sum_{\substack{n \in \mathbb{Z}^{2} \\\|n\| \leqslant N}}\left\langle\varrho_{n}, \mathcal{K}_{T} \varrho_{n}\right\rangle_{\mathcal{H}_{A_{M, c}}}=\sum_{T(x)=x}\left|\operatorname{det}\left(\mathrm{id}-\mathrm{D}_{x} T\right)\right|^{-1} .
$$

Using Lemma 1.4.1, and the definitions (1.1) and (1.9) for the dynamical determinant and Fredholm determinant, respectively, we see directly that

$$
\begin{equation*}
\operatorname{det}\left(1-z \mathcal{K}_{T}\right)=d_{T}(z) \tag{1.24}
\end{equation*}
$$

The Ruelle resonances correspond to the zeroes of the Fredholm determinant, hence to the inverses of the non-zero eigenvalues of $\mathcal{K}_{T}$.

Remark 1.4.2. In view of Equation 1.24 and the relation of the Ruelle resonances of $T$ to the eigenvalues of $\mathcal{K}_{T}$, one may ask how the spectrum of $\mathcal{K}_{T}$ would be affected if we let $\mathcal{K}_{T}$ act on a different Banach space. The following relates a part of the eigenvalues of two linear operators sharing a common dense subspace and is due to a result of Baladi and Tsujii [11, Appendix A]. Consider two separable Banach spaces $\left(\mathcal{B}_{1},\|\cdot\|_{1}\right)$ and $\left(\mathcal{B}_{2},|\cdot|\right)$. This induces two other

Banach spaces

$$
\begin{aligned}
& \quad\left(\mathcal{B}_{1}+\mathcal{B}_{2},\|\cdot\|_{+}\right) \text {and }\left(\mathcal{B}_{1} \cap \mathcal{B}_{2},\|\cdot\|_{\cap}\right) \text {, where } \\
& \|f\|_{+}:=\inf \left\{\left\|f_{1}\right\|_{1}+\left|f_{2}\right| \quad \mid \quad f_{1} \in \mathcal{B}_{1}, f_{2} \in \mathcal{B}_{2}, f=f_{1}+f_{2}\right\} \text { and } \\
& \|f\|_{\cap}:=\max \left\{\|f\|_{1},|f|\right\} \text {. }
\end{aligned}
$$

Suppose that $\mathcal{B}_{\cap}$ is dense in $\mathcal{B}_{1}$ and $\mathcal{B}_{2}$. Let $\mathcal{L}: \mathcal{B}_{+} \rightarrow \mathcal{B}_{+}$be a linear map which preserves the spaces $\mathcal{B}_{\cap}, \mathcal{B}_{1}$ and $\mathcal{B}_{2}$ and is a bounded linear operator on the restrictions $\mathcal{L}_{\mid \mathcal{B}_{1}}$ and $\mathcal{L}_{\mid \mathcal{B}_{2}}$. Then the part of the spectrum of $\mathcal{L}_{\mid \mathcal{B}_{1}}$ and of $\mathcal{L}_{\mid \mathcal{B}_{2}}$ which lies outside the closed disc with radius larger to both essential spectral radii of $\mathcal{L}_{\mid \mathcal{B}_{1}}$ and $\mathcal{L}_{\mid \mathcal{B}_{2}}$ coincide. Moreover, the corresponding generalized eigenspaces of $\mathcal{L}_{\mid \mathcal{B}_{1}}$ and $\mathcal{L}_{\mid \mathcal{B}_{2}}$ coincide and are contained in $\mathcal{B}_{\cap}$.
For the applications that we have in mind, the map $\mathcal{L}$ is just the Koopman or transfer operator, defined on $\mathcal{B}_{1}$ and $\mathcal{B}_{2}$, respectively, extended to the space $\mathcal{B}_{+}$.

The spectrum $\operatorname{sp}\left(\mathcal{K}_{T}\right)$ of $\mathcal{K}_{T}$ on $\mathcal{H}_{A_{M, c}}$ is invariant under complex conjugation since $T$ is real. The constant functions on $\mathbb{T}^{2}$ are all fixed by $\mathcal{K}_{T}$. Therefore we have $1 \in \operatorname{sp}\left(\mathcal{K}_{T}\right)$. If we take $T=M^{k}, k \in \mathbb{N}$ in Lemma 1.4.1, it follows that $\operatorname{tr} \mathcal{K}_{T}=1$. Hence, the dynamical determinant is just $d_{T}(z)=1-z$, also noted in [55, p.3]. We find immediately that 1 is the only Ruelle resonance. We show now that this finding is non-generic in the following sense. The rest of this section is devoted to an idea of Naud [51]. We put for every $r>0$

$$
\begin{equation*}
\mathcal{B}_{r}:=\left\{T \in \mathcal{T}_{r} \mid \text { The lift of } T \text { to } \mathbb{R}^{2} \text { is } \mathbb{Z}^{2} \text {-periodic }\right\} . \tag{1.25}
\end{equation*}
$$

Endowed with the uniform norm this is a Banach space.
Theorem 1.4.3 (Non-trivial Ruelle resonances (I)). Let $M \in \mathrm{SL}_{2}(\mathbb{Z})$ be hyperbolic. For each $r>0$ there exists an open and dense set $\mathcal{G} \subset \mathcal{B}_{r}$ such that the linear functional

$$
B_{M}: \mathcal{B}_{r} \rightarrow \mathbb{R}: \psi \mapsto N_{M}^{-1} \sum_{M x=x} \operatorname{tr}\left((\operatorname{id}-M)^{-1} \mathrm{D}_{x} \psi\right)
$$

never vanishes on $\mathcal{G}$. For all $\psi \in \mathcal{G}$ there exists $\epsilon_{0}>0$ such that for all $\epsilon<\epsilon_{0}$

$$
\operatorname{tr} \mathcal{K}_{M+\epsilon \psi}=1+\epsilon B_{M}(\psi)+O\left(\epsilon^{2}\right) .
$$

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In particular, for all sufficiently small $\epsilon>0$ it holds

$$
\operatorname{sp}\left(\mathcal{K}_{M+\epsilon \psi}\right) \backslash\{0,1\} \neq \varnothing .
$$

Lemma 1.4.4 (Real analyticity of fixed points). Let $M \in \mathrm{SL}_{2}(\mathbb{Z})$ be hyperbolic and $r>0$. Then for all $\psi \in \mathcal{B}_{r}$ the fixed points of the map

$$
M+\delta \psi
$$

are real analytic functions of $\delta$ where $\delta$ lies in a real neighborhood of 0 .
Proof. We set for $\delta \in \mathbb{R}$

$$
F(\delta, x):=M x+\delta \psi(x)-x .
$$

We fix a point $y_{j}:=\left(0, x_{j}\right)$ where $x_{j}, 1 \leqslant j \leqslant N_{M}$, is a fixed point of $M$. By construction, the map $F$ has a holomorphic extension to $\mathbb{C} \times \mathbb{A}_{r}$. Since $M$ is hyperbolic, we have $\operatorname{det} \mathrm{D}_{x_{j}}(F(0, \cdot)) \neq 0$. We apply the Holomorphic Implicit Function Theorem [47, Theorem 1.4.11] on $F$ with $F\left(y_{j}\right)=0$. This yields a holomorphic function $x_{j}(\delta)$ such that $x_{j}(0)=x_{j}$ and which is obviously real analytic for $\delta \in \mathbb{R}$ in a neighborhood of 0 .

Proof of Theorem 1.4.3. Let $\delta \in \mathbb{R}$ and $\psi \in \mathcal{B}_{r}$ and set $\tilde{M}:=M+\delta \psi$. We choose $\delta$ small in Lemma 1.4.4 which gives for each fixed point $x$ of $M$ a real analytic function $\tilde{x}$ with $\tilde{x}(0)=x$. Using a Taylor expansion on $\tilde{x}$ at 0 , we have

$$
\tilde{x}(\delta)=x+O(\delta) .
$$

Using real analyticity of the derivative $\mathrm{D}_{x} \psi$, we have

$$
\mathrm{D}_{x} \psi-\mathrm{D}_{\tilde{x}(\delta)} \psi=O(\delta) .
$$

We write now for each fixed point $x$ of $M$

$$
\begin{aligned}
\operatorname{det}\left(\mathrm{id}-\mathrm{D}_{\tilde{x}(\delta)} \tilde{M}\right) \mid & =\left|\operatorname{det}\left(\mathrm{id}-M-\delta \mathrm{D}_{x} \psi+\delta\left(\mathrm{D}_{x} \psi-\mathrm{D}_{\tilde{x}(\delta)} \psi\right)\right)\right| \\
& =N_{M}\left|\operatorname{det}\left(\mathrm{id}-(\mathrm{id}-M)^{-1}\left(\delta \mathrm{D}_{x} \psi+\left(\delta \mathrm{D}_{x} \psi-\delta \mathrm{D}_{\tilde{x}(\delta)} \psi\right)\right)\right)\right| \\
& =N_{M}\left|\operatorname{det}\left(\mathrm{id}-\delta(\mathrm{id}-M)^{-1} \mathrm{D}_{x} \psi+O\left(\delta^{2}\right)\right)\right| \\
& =N_{M}\left(1-\delta \operatorname{tr}\left((\mathrm{id}-M)^{-1} \mathrm{D}_{x} \psi\right)+O\left(\delta^{2}\right)\right) .
\end{aligned}
$$

We have by Lemma 1.4.1 for $\delta$ small enough

$$
\operatorname{tr} \mathcal{K}_{\tilde{M}}=1+\frac{\delta}{N_{M}} \sum_{M x=x} \operatorname{tr}\left((\mathrm{id}-M)^{-1} \mathrm{D}_{x} \psi\right)+O\left(\delta^{2}\right) .
$$

Now we set

$$
B_{M}: \mathcal{B}_{r} \rightarrow \mathbb{R}: \psi \mapsto N_{M}^{-1} \sum_{M x=x} \operatorname{tr}\left((\mathrm{id}-M)^{-1} \mathrm{D}_{x} \psi\right) .
$$

We next check that this is a non-trivial linear functional. Note that formally $B_{M}(\mathrm{id}-M)=2$. However, no non-zero linear map is in the space of additive perturbations $\mathcal{B}_{r}$. We denote by $v_{j}, j \in\{1,2\}$ the $j$-th column of the matrix $\left((\mathrm{id}-M)^{*}\right)^{-1}$ and we fix now $j$. Let $\psi_{0}: \mathbb{T}+\mathrm{i}(-r, r) \rightarrow \mathbb{C}$ be holomorphic and bounded. For every $\left(x_{1}, x_{2}\right)=: x \in \mathbb{T}^{2}$ we put

$$
\psi(x):=\psi_{0}\left(x_{j}\right) v_{j} .
$$

By construction, we have $\psi \in \mathcal{B}_{r}$ and we evaluate

$$
B_{M}(\psi)=\frac{v_{j}^{*} v_{j}}{N_{M}} \sum_{M x=x} \psi_{0}^{(1)}\left(x_{j}\right)
$$

The right-hand side is a finite sum and by taking for $\psi_{0}$ a suitable Fourier polynomial (e.g. a shifted sine with sufficiently high frequency), we can establish $B_{M}(\psi) \neq 0$. We set $\mathcal{G}:=B_{M}^{-1}(\mathbb{R} \backslash\{0\})$. By continuity of $B_{M}$, the set $\mathcal{G}$ is open and dense in $\mathcal{B}_{r}$.

### 1.5 Non-trivial resonances for the transfer operator

As before, we consider maps $T \in \mathcal{T}_{r}, r>0$ which are sufficiently $C^{1}$-close to a hyperbolic linear map $M \in \mathrm{SL}_{2}(\mathbb{R})$. We turn to the adjoint of $\mathcal{K}_{T}$, acting on the dual space $\mathcal{H}_{A_{M, c}}^{*}$, which we denote by $\mathcal{L}_{T}$.

Lemma 1.5.1 (Transfer operator). Let $M \in \mathrm{SL}_{2}(\mathbb{Z})$ be hyperbolic and let $r>0$. Then there exist constants $0<\delta_{M}$ and $c>0$ such that for each $T \in \mathcal{T}_{r}$ with $d(T, M) \leqslant \delta_{M}$ the map

$$
\mathcal{L}_{T}: \mathcal{H}_{A_{M, c}}^{*} \rightarrow \mathcal{H}_{A_{M, c}}^{*}: f \mapsto \frac{f}{|\operatorname{det} \mathrm{D} T|} \circ T^{-1}
$$

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defines a nuclear operator of order 0 , conjugate to $\mathcal{K}_{T}$. In particular,

$$
\operatorname{sp}\left(\mathcal{L}_{T}\right)=\operatorname{sp}\left(\mathcal{K}_{T}\right) .
$$

Proof. By Theorem 1.3.1 there is $0<\delta_{M}, c>0$ and $\mathcal{H}_{A_{M, c}}$ such that $\mathcal{K}_{T}$ acting on $\mathcal{H}_{A_{M, c}}$ is nuclear of order 0 if $d(T, M) \leqslant \delta_{M}$. The same can be said about its adjoint, acting on $\mathcal{H}_{A_{M, c}}^{*}$ (e.g. see [57, p. 77]). The trace of $\mathcal{K}_{T}$ and $\mathcal{L}_{T}$ coincide, so does their Fredholm determinant, and hence their resonances. By definition of the adjoint, $\forall f^{*} \in \mathcal{H}_{A_{M, c}}^{*}, \forall g \in \mathcal{H}_{A_{M, c}}:\left(\mathcal{L}_{T} f\right)^{*}(g)=f^{*}\left(\mathcal{K}_{T} g\right)$. Using Lemma 1.2.5, it holds

$$
\begin{aligned}
f^{*}\left(\mathcal{K}_{T} g\right) & =\left\langle A_{M, c}^{-2} f, \mathcal{K}_{T} g\right\rangle_{\mathcal{H}_{A_{M, c}}}=\int_{\mathbb{T}^{2}}\left(A_{M, c}^{-1} \bar{f}\right)(x)\left(A_{M, c} \mathcal{K}_{T} g\right)(x) \mathrm{d} x \\
& =\int_{\mathbb{T}^{2}} \bar{f}(x)\left(\mathcal{K}_{T} g\right)(x) \mathrm{d} x=\int_{\mathbb{T}^{2}}\left(\frac{\bar{f}}{|\operatorname{det} \mathrm{D} T|} \circ T^{-1}\right)(x) g(x) \mathrm{d} x \\
& =\left\langle A_{M, c}^{-2}\left(\frac{f}{|\operatorname{det} \mathrm{D} T|} \circ T^{-1}\right), g\right\rangle_{\mathcal{H}_{A_{M, c}}}=\left(\frac{f}{|\operatorname{det} \mathrm{D} T|} \circ T^{-1}\right)^{*}(g) .
\end{aligned}
$$

By Lemma 1.5.1, recalling (1.6), and Lemma 1.4.1 it holds

$$
\operatorname{tr} \mathcal{L}_{T}=\sum_{n \in \mathbb{Z}^{2}} \mathcal{L}_{T} \varrho_{n}^{*}\left(\varrho_{n}\right)=\sum_{T(x)=x}|\operatorname{det}(\mathrm{id}-\mathrm{D} T)|^{-1}
$$

We have the equality

$$
d_{T}(z)=\operatorname{det}\left(1-z \mathcal{K}_{T}\right)=\operatorname{det}\left(1-z \mathcal{L}_{T}\right) .
$$

We give now analogously to Theorem 1.4.3 a spectral result for the transfer operator (recall $\mathcal{B}_{r}$ from (1.25)).

Lemma 1.5.2 (Non-trivial Ruelle resonances (II)). Let $M \in \mathrm{SL}_{2}(\mathbb{Z})$ be hyperbolic. For each $r>0$ there exists an open and dense set $\mathcal{G} \subset \mathcal{B}_{r}$ such that for all $\psi \in \mathcal{G}$ there exists $\epsilon_{0}>0$ such that for all $0<\epsilon \leqslant \epsilon_{0}$

$$
\operatorname{sp}\left(\mathcal{L}_{M+\epsilon \psi}\right) \backslash\{0,1\} \neq \varnothing .
$$

Proof. By Theorem 1.4.3 we know that under every perturbation $\psi \in \mathcal{G}$ there is $\epsilon_{0}>0$ such that we find for all $0<\epsilon \leqslant \epsilon_{0}$ non-trivial Ruelle resonances. Using Lemma 1.5 .1 for well-definedness of $\mathcal{L}_{M+\epsilon \psi}$ and for the relation $\operatorname{sp}\left(\mathcal{L}_{T}\right)=$ $\operatorname{sp}\left(\mathcal{K}_{T}\right)$, the result follows.

Clearly, the Lebesgue measure (by Remark 1.2.6, the constant density 1 ) is fixed by $\mathcal{L}_{M}$. This does not persist under a generic perturbation of $M$. However, the spectral relation in Lemma 1.5 .1 implies that $\mathcal{L}_{T}$ fixes some functionals in $\mathcal{H}_{A_{M, c}}^{*}$. In particular, using Remark 1.4.2, we can apply [15, Theorem 3] to our transfer operators $\mathcal{L}_{M}$ and $\mathcal{L}_{T}$. Hence, the eigenvalue 1 of $\mathcal{L}_{T}$ is simple and the projector $\Pi_{1}^{*}$ onto the corresponding eigenspace of $\mathcal{L}_{T}$ gives us the SRB measure

$$
\mu_{\mathrm{SRB}}:=\Pi_{1}^{*} 1^{*},
$$

in the usual sense. (It is absolutely continuous with respect to Lebesgue measure in the unstable direction.)
We finish this section by showing the existence of non-zero perturbations $\psi \in \mathcal{B}_{r}$ which allow the determinant $\operatorname{det}\left(M+\epsilon \mathrm{D}_{x} \psi\right)$ to remain constant or to vary for $x \in \mathbb{T}^{2}$.

Lemma 1.5.3 (Volume under perturbations). Let $r>0$ and let $M \in \mathrm{SL}_{2}(\mathbb{Z})$ be hyperbolic. Then there exist non-zero maps $\psi \in \mathcal{B}_{r}$ in each of the following cases:
(i) For all $\epsilon>0$ and all $x \in \mathbb{T}^{2}$ it holds $\operatorname{det}\left(M+\epsilon \mathrm{D}_{x} \psi\right)=1$.
(ii) For all $\epsilon>0$ and Lebesgue almost all $x \in \mathbb{T}^{2}$ it holds $\left|\operatorname{det}\left(M+\epsilon \mathrm{D}_{x} \psi\right)\right| \neq$ 1.

In particular, the map $\psi$ can be chosen such that for all small $\epsilon>0$ the corresponding transfer operator

$$
\mathcal{L}_{M+\epsilon \psi}
$$

admits non-trivial Ruelle resonances.
Proof. We prove first Claim (i), including the statement about the non-trivial Ruelle resonances. We will apply Lemma A. 2 (i). We choose $j \in\{1,2\}, r>0$ and let $\phi: \mathbb{T}+\mathrm{i}(-r, r) \rightarrow \mathbb{C}$ be a holomorphic and bounded map. For $\alpha \in \mathbb{R}^{2}$ we set for every $\left(x_{1}, x_{2}\right)=: x \in \mathbb{T}^{2}$

$$
\psi_{\phi, \alpha}(x):=\left(\alpha_{1} \phi\left(x_{j}\right), \alpha_{2} \phi\left(x_{j}\right)\right)
$$

We put $d:=2, j, T:=M, \phi$ and $T_{\phi}:=\psi_{\phi, \alpha}$ (e.g. as lift to $\mathbb{R}^{2}$ ) in Lemma A.2. Since $M$ is a constant matrix, say, $M=\left(\begin{array}{ll}a & b \\ c & d\end{array}\right)$ for suitable $a, b, c, d \in \mathbb{Z}$, we can write Condition A. 2 (i) as

$$
\begin{equation*}
\alpha_{1} d=\alpha_{2} b \quad \text { if } j=1 \quad \text { or } \quad \alpha_{1} c=\alpha_{2} a \quad \text { if } j=2 \tag{1.26}
\end{equation*}
$$

Hence, we have non-zero solutions in $\alpha$ independent of $x$. We choose such a solution $\alpha$ and take $\psi=\psi_{\phi, \alpha}$. Then $\psi \in \mathcal{B}_{r}$ which yields $\operatorname{det}\left(M+\epsilon \mathrm{D}_{x} \psi\right)=1$ for every $\epsilon>0$. We are free to choose any suitable $\phi$. In particular, Theorem 1.4.3 yields a linear functional $B_{M}$ and a dense subset $\mathcal{G} \subset \mathcal{B}_{r}$ on which $B_{M}$ is non-zero. We have to make sure that $\psi \in \mathcal{G}$. Then for $\epsilon$ small $\mathcal{L}_{M+\epsilon \psi}$ admits non-trivial Ruelle resonances by Lemma 1.5.2. To this end, we evaluate $B_{M}$ at $\psi$ which yields

$$
B_{M}(\psi)=B_{M}\left(\psi_{\phi, \alpha}\right)=N_{M}^{-1} \sum_{M x=x} \operatorname{tr}\left((\operatorname{id}-M)^{-1} \mathrm{D}_{x} \psi_{\phi, \alpha}\right)=\frac{v_{j}^{*} \alpha}{N_{M}} \sum_{M x=x} \phi^{(1)}\left(x_{j}\right),
$$

where $v_{j}^{*}$ is the $j$-th row of $(\mathrm{id}-M)^{-1}$. The sum over the fixed points of $M$ can be made non-zero by a suitable Fourier polynomial. Now we have

$$
v_{1}^{*} \alpha=\frac{(1-d) \alpha_{1}+c \alpha_{2}}{\operatorname{det}(\operatorname{id}-M)} \quad \text { or } \quad v_{2}^{*} \alpha=\frac{b \alpha_{1}+(1-a) \alpha_{2}}{\operatorname{det}(\operatorname{id}-M)} .
$$

Using (1.26), we find

$$
v_{1}^{*} \alpha=\frac{\left(c-b+\frac{b}{d}\right) \alpha_{2}}{\operatorname{det}(\operatorname{id}-M)} \quad \text { or } \quad v_{2}^{*} \alpha=\frac{\left(b-c+\frac{c}{a}\right) \alpha_{1}}{\operatorname{det}(\operatorname{id}-M)} .
$$

Both equations can never be zero since $M$ is not diagonal. We prove now Claim (ii) by modifying the map $\psi$. For $\delta \in \mathbb{R} \backslash\{0\}$ we set $\tilde{\alpha}:=\alpha+\delta w_{j}$, where $w_{j}$ is the $j$-th column of $M$ and put $\tilde{\psi}:=\psi_{\phi, \tilde{\alpha}}$. We have

$$
\operatorname{det}\left(M+\epsilon \mathrm{D}_{x} \tilde{\psi}\right)=\operatorname{det}\left(M+\epsilon \mathrm{D}_{x} \psi+\epsilon \mathrm{D}_{x}(\tilde{\psi}-\psi)\right)=1+\delta \epsilon \phi^{(1)}\left(x_{j}\right)
$$

Since $\phi$ is not constant, the right-hand side differs from 1 (and -1 ) for Lebesgue almost all $x$. Since $v_{j}^{*} \tilde{\alpha}=v_{j}^{*} \alpha+\delta v_{j}^{*} w_{j} \neq 0$ for the right choice of the sign of $\delta$, we have $B_{M}(\tilde{\psi}) \neq 0$.

## 2 Horocycle averages on closed manifolds

### 2.1 Introduction

Let $M$ be a closed (compact without boundary) orientable Riemannian manifold of arbitrary finite dimension $d \geqslant 3$. On such manifolds Anosov introduced $C^{2}$ flows

$$
g_{\alpha_{1}} \circ g_{\alpha_{2}}=g_{\alpha_{1}+\alpha_{2}}=g_{\alpha}: M \rightarrow M, \quad \alpha, \alpha_{1}, \alpha_{2} \in \mathbb{R},
$$

to study the geodesic flow on the unit tangent bundle of closed Riemannian manifolds with variable negative sectional curvature [3], [4]. As pointed out by Anosov [3], the topological entropy $h_{\text {top }}$ of the time-one map $g_{1}$ of an Anosov flow is positive.
A special class of such Anosov flows are those which preserve a contact structure. The geodesic flows are well-studied examples of contact Anosov flows. We give the precise definition of a (contact) Anosov flows in Section 2.2.
Every Anosov flow admits a contracting transversal foliation. The underlying vector bundle $E_{-}$is called the strong stable distribution. If the leaves of the contracting foliation are one-dimensional and orientable, one associates with $g_{\alpha}$ another flow, the horocycle flow $h_{\rho}: M \rightarrow M, \rho \in \mathbb{R}$. (The term horocycle flow was used originally only in the case of the geodesic flow, e.g. see [49, p.84] or [38].) For every $x \in M$ the flow trajectory $h_{\mathbb{R}}(x)$ is such a contracting leaf. Statistical properties of contact Anosov flows are nowadays fairly well understood (see [23], [33], [48]). Regarding the horocycle flow one knows by the work of Bowen and Marcus unique ergodicity of and minimality of the horocycle flow (e.g. see [18], [50]). The corresponding invariant probability measure $\mu$ will play an important role below. (It is related to but distinct from the measure of maximal entropy of the flow.)
Since the horocycle flow is induced by the Anosov flow the following pointwise

2 Horocycle averages on closed manifolds
equality for all $x \in M$ holds for a suitable function $\tau(\rho, \alpha, x)$ :

$$
g_{\alpha} \circ h_{\rho}(x)=h_{\tau(\rho, \alpha, x)} \circ g_{\alpha}(x) .
$$

We call $\tau(\rho, \alpha, x)$ the renormalization time.
This kind of renormalization has been used effectively in the work of Flaminio and Forni [28] to give a precise understanding of the horocycle integral

$$
\gamma_{x}(\varphi, T):=\int_{0}^{T} \varphi \circ h_{\rho}(x) \mathrm{d} \rho, \quad x \in M, \quad T>0,
$$

in the setting of unit speed geodesic flows on hyperbolic compact (more generally finite volume) Riemannian surfaces with constant negative sectional curvature (i.e. Riemann surfaces), for $\varphi: M \rightarrow \mathbb{R}$ in Sobolev spaces of positive order. In this case, $h_{\text {top }}=1$. Flaminio and Forni found that the speed of convergence of $\gamma_{x}(\varphi, T) / T$ to $\mu(\varphi)$ as $T \rightarrow \infty$ is controlled by invariant distributions under the push-forward of the horocyclic vector field. These distributions are also eigendistributions under the push-forward of the geodesic vector field and the eigenvalues give the powers of $T$ appearing in the expansion of $T^{-1} \gamma_{x}(\varphi, T)-\mu(\varphi)$.
Their approach inspired Giulietti and Liverani [30] to study a toy model, replacing the Anosov flow with a hyperbolic diffeomorphism, using the renormalization dynamics as a key to study $\gamma_{x}(\varphi, T)$. They show analogously (for the corresponding invariant measure $\mu$ ) that the speed of convergence to zero of $T^{-1} \gamma_{x}(\varphi, T)-\mu(\varphi)$ is controlled by eigendistributions for a weighted transfer operator of the hyperbolic diffeomorphism.
Giulietti and Liverani conjectured that a similar behavior holds in the setting of more general Anosov flows, e.g. for the geodesic flow on the unit tangent bundle of a Riemannian manifold with variable negative sectional curvature [30, Conjecture 2.12]. More precisely, we expect for smooth enough observables $\varphi$ an expansion like

$$
\begin{equation*}
\gamma_{x}(\varphi, T)=T \int \varphi \mathrm{~d} \mu+\sum_{\delta<\Re \lambda<h_{\text {top }}} T^{\frac{\Re \lambda}{h_{\text {top }}}} c(\lambda, T, x) \mathcal{O}_{\lambda}(\varphi)+\mathcal{E}_{T, x}(\varphi), \tag{2.1}
\end{equation*}
$$

with $\mathcal{E}_{T, x}=O\left(T^{\frac{\delta}{h_{\text {top }}}}\right)$, uniformly in $x$. The $\mathcal{O}_{\lambda}$ are generalized eigendistributions associated to the eigenvalue $\lambda$ for the adjoint of the generator $X+V$ of a certain weighted transfer operator $\mathcal{L}_{\alpha, \phi_{\alpha}}$, acting on an anisotropic Banach space (see below). The real parameter $\delta$ is an upper bound on the essential spectral bound
of $X+V$. The complex coefficients $c(\lambda, T, x)$ are bounded from above independently of $x$ by $|\log T|^{c}$ for some $c=c(\lambda) \geqslant 0$ which depends whether $\Re \lambda<0$, $\Re \lambda=0$ or $\Re \lambda>0$ and if there are non-trivial Jordan blocks for $\lambda$. This is analogous to the bounds in [28],[30]. However our methods show no substantial improvement of the error term $\mathcal{E}_{T, x}$ if the summation in $\lambda$ includes some $\Re \lambda<0$ (this is seen also in [28],[30]). We restrict ourself therefore to $\delta \geqslant 0$ (i.e. always $\Re \lambda>0)$.

The main result of this work, Theorem 2.5.7, gives conditions under which such an asymptotic expansion indeed holds, for some $\delta>0$, for codimension one topologically mixing Anosov flows, under an assumption of "spectral gap with (Dolgopyat) bounds" (Condition 2.4.11 below). In Proposition 2.5.10 we specialize to $C^{3}$ contact Anosov flows in dimension $d=3$. For compact Riemann surfaces (recall that this is the constant negative curvature case) Randol [53] proved that there exist eigenvalues arbitrarily close to 1 (his result is for the associated Laplacian). This provides examples with a non-trivial expansion. Analogous to the work of Giulietti and Liverani [30], the key idea to study $\gamma_{x}(\varphi, T)$ is to introduce a weighted transfer operator family

$$
\mathcal{L}_{\alpha, \phi_{\alpha}}: W_{p}^{s, t, q} \rightarrow W_{p}^{s, t, q}, \quad \mathcal{L}_{\alpha, \phi_{\alpha}} \varphi=\phi_{\alpha} \cdot \varphi \circ g_{-\alpha}, \quad \alpha \geqslant 0
$$

where the weight is $\phi_{\alpha}=\partial_{\rho} \tau(0,-\alpha, \cdot)$ and where $W_{p}^{s, t, q}$ is an anisotropic Banach space with certain real regularity parameters $s, t, q$ and $p$. In the case of the unit speed parametrization of the flow $h_{\rho}$, the weight $\partial_{\rho} \tau(0,-\alpha, \cdot)$ is just the Jacobian along the strong stable distribution evaluated at negative time $-\alpha$.

The paper is organized as follows: After recalling some facts about Anosov flows in Section 2.2, the transfer operator $\mathcal{L}_{\alpha, \phi_{\alpha}}$ is defined in Section 2.3.1 (for more general weights) and the Banach spaces $W_{p}^{s, t, q}$ are constructed in Section 2.3.2. These spaces are a flow analogue to the spaces constructed by Baladi and Tsujii [10] to study hyperbolic diffeomorphisms. Anisotropic Banach spaces are now considered a standard tool (yet with still ongoing research) for investigating transfer operators and zeta functions associated to hyperbolic dynamics [7]-[9], [13], [15], [23], [31], [33], [48], [52], [64], [65]. Although we do not study here the dynamical zeta function for the transfer operator $\mathcal{L}_{\alpha, \phi_{\alpha}}$, we believe that this space could be a suitable choice to be dealt with.
In Section 2.4 we establish properties of the transfer operator, its generator $X+V$ and the resolvent $\mathcal{R}_{z}$. Most of these results do not require the contact
assumption. Among those are norm estimates which yield a Lasota-York inequality for the resolvent. This is Theorem 2.4.5. Then in Lemma 2.4.10 one obtains a strip in the spectrum of the generator, containing at most countable eigenvalues of finite multiplicity. Those are precisely the eigenvalues $\lambda$ in the summation over $\lambda$ in (2.1). Finally, these results are used in Section 2.5 to give the expansion (2.1) of $\gamma_{x}(\varphi, T)$ in terms of eigendistributions and eigenvalues of $X+V$ under a spectral gap with bounds condition, see Condition 2.4.11.

We end this introduction with two remarks about possible further work:
First, the conjecture that the distributions $\mathcal{O}_{v}$ appearing in the expansion (2.1) are fixed by the (adjoint) of the horocycle flow remains still open. (In contrast this was the starting point in [28]!) Here, progress has been made by Faure and Guillarmou [24] in dimension 3 for smooth contact Anosov flows.

Second, the renormalization time $\tau(\rho, \alpha, x)$ inherits the regularity properties of the underlying Anosov foliation and horocycle flow, i.e. the regularity in $x$ is expected to be no more than Hölder. To deal with such irregular flows one can lift the dynamics to the Grassmanian. This has been used with success, e.g. in [30], [33] and more recently in [65]. However in this work we wish to avoid such technicalities and we will make additional assumptions ensuring that $\tau(\rho, \alpha, x)$ enjoys sufficient regularity.
In particular, if the Anosov flow is $C^{r}$ we require $\partial_{\rho} \tau(0, \alpha, \cdot)$ to be $C^{r-1}$ for all $\alpha \geqslant 0$. This is reasonably only if $r$ is small since the regularity of the stable foliation is usually only Hölder. In the setting of $C^{3}$ contact Anosov flows in dimension 3 we can take $r=2-\epsilon$ for all $\epsilon>0$ by a result of [40] (see also Remark 2.5.8 in Section 2.5).

The Appendix comprises our computational tools. On the lowest level, we utilize Fourier transform, integration by parts, and Young's inequality [17, Theorem 3.9.4] to estimate convolutions.

### 2.2 Geometric setting

Let $M$ be a closed, connected, orientable, smooth Riemannian manifold of dimension $d \geqslant 3$. We let $g_{\alpha}: M \rightarrow M, \alpha \in \mathbb{R}$, be a $C^{r}$ Anosov flow on $M^{1}$ for $r>1$. That is, there exists a decomposition of the tangent space $T M$ of $M$ as

[^1]a direct sum
\[

$$
\begin{equation*}
\mathrm{T} M=E_{-} \oplus E_{+} \oplus E_{0} \tag{2.2}
\end{equation*}
$$

\]

such that for some constants $C \geqslant 1,0<\theta<1$ and every $\alpha \geqslant 0$

$$
\begin{align*}
&\left\|\mathrm{D} g_{\alpha} v\right\| \leqslant \theta^{\alpha}\|v\|,  \tag{2.3}\\
&\left\|\mathrm{D} g_{-\alpha} v\right\| \text { for all } v \in E_{-}, \\
& \\
&,
\end{align*}
$$

and $E_{0}=\langle X\rangle$ where $X$ is the generator of the Anosov flow

$$
\begin{equation*}
X:=\partial_{\alpha} g_{-\alpha \mid \alpha=0} . \tag{2.4}
\end{equation*}
$$

Note that the conditions in (2.3) are closed. Hence by compactness of $M$ the distributions $E_{-}$and $E_{+}$are uniformly continuous and so are the weak-stable $E_{-} \oplus E_{0}$ and weak-unstable $E_{+} \oplus E_{0}$ distributions. The restriction of the tangent space to a base point $x \in M$ is denoted by

$$
\begin{equation*}
\mathrm{T}_{x} M=E_{-, x} \oplus E_{+, x} \oplus E_{0, x}, \tag{2.5}
\end{equation*}
$$

The dimensions of those vector spaces do not vary with $x$ and we set for some $x \in M$

$$
\begin{equation*}
d_{-}:=\operatorname{dim} E_{-, x} . \tag{2.6}
\end{equation*}
$$

The cotangent space $T^{*} M$ is the dual space of $T M$ and has the canonical splitting

$$
\begin{equation*}
\mathrm{T}^{*} M=E_{-}^{*} \oplus E_{+}^{*} \oplus E_{0}^{*} \quad \text { and } \quad \mathrm{T}_{x}^{*} M=E_{-, x}^{*} \oplus E_{+, x}^{*} \oplus E_{0, x}^{*}, \quad x \in M, \tag{2.7}
\end{equation*}
$$

where $E_{-}^{*} \cong\left(E_{+} \oplus E_{0}\right)^{\perp}, E_{+}^{*} \cong\left(E_{-} \oplus E_{0}\right)^{\perp}, E_{0}^{*} \cong\left(E_{-} \oplus E_{+}\right)^{\perp}$. This splitting is $\left(\mathrm{D} g_{\alpha}\right)^{\mathrm{tr}}$-invariant and satisfies an analogue of (2.3).
A contact form is a 1 -form $\eta \in T^{*} M$ such that $\eta \wedge \bigwedge_{n=1}^{\frac{d-1}{2}} \mathrm{~d} \eta$ vanishes nowhere ( $\mathrm{d} \eta$ is the exterior derivative of $\eta$ ). An Anosov flow is a contact flow if there exists a $C^{1}$ contact form $\eta$ which is preserved by the pullback of $g_{\alpha}$. Clearly, a contact form can only exist if $d$ is odd.
We mean by " $\Subset$ " for sets $A, B \subseteq T^{*} M$ (or $\subseteq \mathbb{R}^{d}$ ) that

$$
A \Subset B \Leftrightarrow \bar{A} \subseteq(\operatorname{int} B \cup\{0\}) .
$$

Here $\bar{A}$ denotes the closure of $A$ and $\operatorname{int} B$ the interior of $B$. We say that a cone $A$ is compactly included in a cone $B$ if and only if $A \Subset B$. We say that a cone $A$ and a cone $B$ are transversal if and only if $A \cap B=\{0\}$.
We introduce two closed convex cone fields on $M$ in the cotangent space:
For every $x \in M$ and for every $v \in \mathrm{~T}_{x}^{*} M$ we have $v=v^{-}+v^{+}+v^{0}$, where $v^{\sigma} \in E_{\sigma, x}^{*}, \sigma \in\{-,+, 0\}$. For every $0<\gamma<1$ we set

$$
\begin{align*}
C_{\gamma}^{-}(x) & :=\left\{v \in \mathrm{~T}_{x}^{*} M \mid\left\|v^{+}\right\|+\left\|v^{0}\right\| \leqslant \gamma\left\|v^{-}\right\|\right\},  \tag{2.8}\\
C_{\gamma}^{+}(x) & :=\left\{v \in \mathrm{~T}_{x}^{*} M \mid\left\|v^{-}\right\|+\left\|v^{0}\right\| \leqslant \gamma\left\|v^{+}\right\|\right\} .
\end{align*}
$$

If $\gamma^{\prime}>\gamma$ then we have the compact inclusions

$$
C_{\gamma}^{-}(x) \Subset C_{\gamma^{\prime}}^{-}(x) \quad \text { and } \quad C_{\gamma}^{+}(x) \Subset C_{\gamma^{\prime}}^{+}(x)
$$

Moreover, this construction implies $E_{-, x}^{*} \subset C_{\gamma}^{-}(x)$, and $E_{+, x}^{*} \subset C_{\gamma}^{+}(x)$ and also transversality $E_{0, x}^{*} \cap\left(C_{\gamma}^{-}(x) \cup C_{\gamma}^{+}(x)\right)=\{0\}$ and $C_{\gamma}^{-}(x) \cap C_{\gamma}^{+}(x)=\{0\}$.
We have (see Lemma B.1) for all $\alpha \geqslant 0$ so that $C^{2} \theta^{\alpha} \gamma<\gamma^{\prime}<1$ and for all $x \in M$ the compact inclusions

$$
\begin{equation*}
\left(\mathrm{D} g_{-\alpha}\right)^{\operatorname{tr}} C_{\gamma}^{-}(x) \Subset C_{\gamma^{\prime}}^{-}\left(g_{\alpha}(x)\right) \quad \text { and } \quad\left(\mathrm{D} g_{\alpha}\right)^{\operatorname{tr}} C_{\gamma}^{+}(x) \Subset C_{\gamma^{\prime}}^{+}\left(g_{-\alpha}(x)\right) . \tag{2.9}
\end{equation*}
$$

The cones defined in (2.8) are expanding and contracting, respectively (see Lemma B.2). Note that the cones in (2.8) have non-empty interior while [43, Proposition 17.4.4] uses "flat" cones included in $E_{+}^{*} \oplus E_{-}^{*}$.
Let $V_{\omega} \subseteq \Omega, \omega \in \Omega$, be an open cover of $M$, where $\Omega$ is a finite index set. We let $\mathcal{A}$ be an atlas for $M$, containing diffeomorphic $C^{r}$-charts $\kappa_{\omega}: V_{\omega} \rightarrow \mathbb{R}^{d}$, compatible with the splittings (2.2) and (2.7), as we explain now. Fixing coordinates $\left(x_{1}, \ldots, x_{d}\right) \in \mathbb{R}^{d}$ and recalling $X$ from (2.4), we may and do require the flowbox condition

$$
\begin{equation*}
\mathrm{D} \kappa_{\omega} X_{\mid V_{\omega}}=\partial_{x_{d}} . \tag{2.10}
\end{equation*}
$$

Since $g_{\alpha}$ is $C^{r}$ the chart maps $\kappa_{\omega}, \omega \in \Omega$, are also $C^{r}$ diffeomorphisms. We set

$$
\begin{equation*}
C_{\gamma, \omega}^{\sigma}:=\bigcup_{x \in V_{\omega}}\left(\mathrm{D} \kappa_{\omega}^{-1}\right)^{\operatorname{tr}} C_{\gamma}^{\sigma}(x), \quad \sigma \in\{-,+\}, \omega \in \Omega . \tag{2.11}
\end{equation*}
$$

We require the sets $V_{\omega}$ to be small enough such that for small $0<\gamma_{-}, \gamma_{+} \leqslant 1$ there exist $0<\gamma_{-}^{*}, \gamma_{+}^{*} \leqslant 1$ such that for all $\omega \in \Omega$ and for all $x \in V_{\omega}$

$$
\begin{equation*}
\left(\mathrm{D}_{x} \kappa_{\omega}\right)^{\operatorname{tr}} C_{\gamma_{-}, \omega}^{-} \subseteq C_{\gamma_{-}^{*}}^{-}(x) \quad \text { and } \quad\left(\mathrm{D}_{x} \kappa_{\omega}\right)^{\operatorname{tr}} C_{\gamma_{+}, \omega}^{+} \subseteq C_{\gamma_{+}^{* *}}^{+}(x) . \tag{2.12}
\end{equation*}
$$

This is possible by uniform continuity of the weak-stable and weak-unstable distributions and the flowbox condition in (2.10). Note that the cones $C_{\gamma, \omega}^{\sigma}$ are not necessarily convex. This poses no problem since the differential is linear and hence the convex closure of $C_{\gamma_{\sigma}, \omega}^{\sigma}$ is contained in $C_{\gamma_{\sigma}^{*}}^{\sigma}(x)$ (this is already a convex, closed cone) for all $x \in M$. Without loss of generality we identify $C_{\gamma_{\sigma}, \omega}^{\sigma}$ with its convex closure.

Definition 2.2.1 (Cone ensemble). Let $C^{-}, C^{+} \subset \mathbb{R}^{d}, d \geqslant 3$, be transversal, convex, closed cones with non-empty interiors. Let $\Phi_{\sigma}: \mathbb{R}^{d} \backslash\{0\} \rightarrow[0,1]$ be $C^{\infty}$ maps, $\sigma \in\{-,+, 0\}$, such that

$$
\begin{aligned}
& \Phi_{-\mid \operatorname{int} C^{-}} \equiv 1, \quad \Phi_{+\mid \operatorname{int} C^{+}} \equiv 1, \quad \Phi_{-}+\Phi_{+}+\Phi_{0} \equiv 1 \quad \text { and } \\
& C^{-}=\mathbb{R}^{d} \backslash\left(\operatorname{supp} \Phi_{+} \cup \operatorname{supp} \Phi_{0}\right), C^{+}=\mathbb{R}^{d} \backslash\left(\operatorname{supp} \Phi_{-} \cup \operatorname{supp} \Phi_{0}\right) .
\end{aligned}
$$

We call $\Theta:=\left(\Phi_{-}, \Phi_{+}, \Phi_{0}\right)$ a cone ensemble. ${ }^{2}$
Definition 2.2.2 (Cone hyperbolicity). Let $K \subset \mathbb{R}^{d}$ be open and let $F: K \rightarrow$ $F(K)$ be a diffeomorphism. Let $\Theta, \Theta^{\circ}$ be two cone ensembles. Let

$$
C^{-}:=\mathbb{R}^{d} \backslash\left(\operatorname{supp} \Phi_{+} \cup \operatorname{supp} \Phi_{0}\right) .
$$

We say that $F$ is $\left(\Theta^{\circ}, \Theta\right)$-cone hyperbolic on $K$ if there exists $C^{\infty}$ maps

$$
\tilde{\Phi}_{+}, \widetilde{\Phi}_{\sigma}^{\circ}: \mathbb{R}^{d} \backslash\{0\} \rightarrow[0,1]
$$

such that $\widetilde{\Phi}_{+\mid \operatorname{supp} \Phi_{+}}, \widetilde{\Phi}_{\sigma \mid \operatorname{supp} \Phi_{\sigma}^{\circ}}^{\circ} \equiv 1$ for all $\sigma \in\{-, 0\}$ such that for all $z \in K$

$$
\begin{equation*}
\left(\mathrm{D}_{z} F\right)^{\operatorname{tr}} \operatorname{supp} \tilde{\Phi}_{-}^{\circ} \Subset C^{-} \text {and }\left(\mathrm{D}_{z} F\right)^{\operatorname{tr}} \operatorname{supp} \widetilde{\Phi}_{0}^{\circ} \Subset \mathbb{R}^{d} \backslash \operatorname{supp} \tilde{\Phi}_{+} . \tag{2.13}
\end{equation*}
$$

In Section 2.3.2 an anisotropic Banach space is constructed where the cones $C^{-}$, $C^{+}$determine the directions of lowest and highest regularity, respectively. The inclusions (2.13) ensure that no parts of higher regularity are mapped to parts of lower regularity.

[^2]2 Horocycle averages on closed manifolds

Lemma 2.2.3 (Existence of admissible cones). Let $\alpha \in \mathbb{R}$ and let $\omega, \omega^{\prime} \in \Omega$. Set $V_{\alpha, \omega \omega^{\prime}}:=V_{\omega} \cap g_{\alpha}\left(V_{\omega^{\prime}}\right)$ and set

$$
F_{-\alpha, \omega \omega^{\prime}}: \kappa_{\omega}\left(V_{\alpha, \omega \omega^{\prime}}\right) \rightarrow \kappa_{\omega^{\prime}}\left(V_{-\alpha, \omega^{\prime} \omega}\right): y \mapsto \kappa_{\omega^{\prime}} \circ g_{-\alpha} \circ \kappa_{\omega}^{-1}(y) .
$$

Then there exists $\alpha_{0}>0$ such that for all $\omega, \omega^{\prime} \in \Omega$ there exist cone ensembles

$$
\Theta_{\omega}=\left(\Phi_{-, \omega}, \Phi_{+, \omega}, \Phi_{0, \omega}\right) \quad \text { and } \quad \Theta_{\omega^{\prime}}^{\circ}=\left(\Phi_{-, \omega^{\prime}}^{\circ}, \Phi_{+, \omega^{\prime}}^{\circ}, \Phi_{0, \omega^{\prime}}^{\circ}\right),
$$

such that for all $\alpha \geqslant \alpha_{0}$ the map $F_{-\alpha, \omega \omega^{\prime}}$ is $\left(\Theta_{\omega^{\prime}}^{\circ}, \Theta_{\omega}\right)$-cone hyperbolic. Moreover, for every $\omega \in \Omega$ it holds

$$
\begin{equation*}
\operatorname{supp} \Phi_{0, \omega}^{\circ} \Subset \operatorname{supp} \Phi_{0, \omega} \cup \operatorname{supp} \Phi_{+, \omega} \quad \text { and } \quad \operatorname{supp} \Phi_{+, \omega}^{\circ} \Subset \operatorname{supp} \Phi_{+, \omega} \tag{2.14}
\end{equation*}
$$

Proof. We let $\omega, \omega^{\prime} \in \Omega$. We assume $V_{\alpha, \omega \omega^{\prime}} \neq \varnothing$ (otherwise we are done). We let $0<\gamma_{-}, \gamma_{+} \leqslant 1$ be small such that $\gamma_{-}^{*}, \gamma_{+}^{*}>0$ are the values attained in (2.12) for all cones $C_{\gamma-, \omega}^{-}, C_{\gamma_{+}, \omega}^{+}, \omega \in \Omega$. These cones are transversal, convex and closed by construction. We repeat the construction, resulting in values $\widetilde{\gamma}_{-}^{*}<\gamma_{-}^{*}$ and $\widetilde{\gamma}_{+}^{*}<\gamma_{+}^{*}$, using now values

$$
\tilde{\gamma}_{-}<\gamma_{-}, \quad \tilde{\gamma}_{+}<\gamma_{+},
$$

sufficiently small (possibly by passing to a finer open cover) such that for all $\omega \in \Omega$ and all $x \in V_{\omega}$

$$
\begin{equation*}
\left(\mathrm{D}_{\kappa_{\omega}(x)} \kappa_{\omega}^{-1}\right)^{\operatorname{tr}} C_{\tilde{\gamma}_{-}^{*}}^{-}(x) \Subset C_{\gamma_{-}, \omega}^{-}, \quad\left(\mathrm{D}_{\kappa_{\omega}(x)} \kappa_{\omega}^{-1}\right)^{\operatorname{tr}} C_{\tilde{\gamma}_{+}^{*}}^{+}(x) \Subset C_{\gamma_{+}, \omega}^{+} . \tag{2.15}
\end{equation*}
$$

We note that the map $F_{\alpha, \omega \omega^{\prime}}$ is a diffeomorphism by construction. We construct further cones as follows: By the construction of local cones in (2.11) and the compact inclusion given in (2.12) for some $C^{2} \beta^{\alpha} \gamma_{+}^{*} \leqslant \gamma_{+}^{\prime}<\tilde{\gamma}_{+}^{*}$ and for all $\alpha \geqslant \alpha_{0}$ we have for all $x \in V_{\alpha, \omega \omega^{\prime}}$

$$
\left(\mathrm{D} g_{\alpha}\right)^{\operatorname{tr}}\left(\mathrm{D}_{x} \kappa_{\omega}\right)^{\operatorname{tr}} C_{\gamma_{+}, \omega}^{+} \subseteq\left(\mathrm{D} g_{\alpha}\right)^{\operatorname{tr}} C_{\gamma_{+}^{*}}^{+}(x) \Subset C_{\gamma_{+}^{\prime}}^{+}\left(g_{-\alpha}(x)\right) \Subset C_{\tilde{\gamma}_{+}^{*}}^{+}\left(g_{-\alpha}(x)\right) .
$$

Comparing with the compact inclusion in (2.15), there exists a convex, closed cone $\widetilde{C}_{\gamma_{+}, \omega^{\prime}}^{+} \subset \mathbb{R}^{d}$ such that

$$
\begin{equation*}
\left(\mathrm{D} F_{\alpha, \omega^{\prime} \omega}\right)^{\operatorname{tr}} C_{\gamma_{+}, \omega}^{+} \Subset \widetilde{C}_{\gamma_{+}, \omega^{\prime}}^{+} \Subset C_{\gamma_{+}, \omega^{\prime}}^{+} \tag{2.16}
\end{equation*}
$$

Analogously we find

$$
\begin{equation*}
\left(\mathrm{D} F_{-\alpha, \omega^{\prime} \omega}\right)^{\operatorname{tr}} C_{\gamma_{-}, \omega}^{-} \Subset \widetilde{C}_{\gamma_{-}, \omega^{\prime}}^{-} \Subset C_{\gamma_{-}, \omega^{\prime}}^{-} \tag{2.17}
\end{equation*}
$$

Recalling Definition 2.2.1, we let

$$
\Theta_{\omega}=\left(\Phi_{-, \omega}, \Phi_{+, \omega}, \Phi_{0, \omega}\right) \quad \text { and } \quad \Theta_{\omega^{\prime}}^{\circ}=\left(\Phi_{-, \omega^{\prime}}^{\circ}, \Phi_{+, \omega^{\prime}}^{\circ}, \Phi_{0, \omega^{\prime}}^{\circ}\right)
$$

be the cone ensembles such that

$$
\Phi_{-, \omega \mid \operatorname{int} \widetilde{C}_{\gamma_{-}, \omega}^{-}} \equiv \Phi_{+, \omega \mid \operatorname{int} C_{\gamma_{+}, \omega}^{+}} \equiv \Phi_{-, \omega^{\prime} \mid \operatorname{int} C_{\gamma_{-}, \omega^{\prime}}^{-}}^{\circ} \equiv \Phi_{+, \omega^{\prime} \mid \operatorname{int} \widetilde{C}_{\gamma_{+}, \omega^{\prime}}^{+}}^{\circ} \equiv 1
$$

The supports of $\Phi_{-, \omega}, \Phi_{+, \omega}$ and $\Phi_{-, \omega^{\prime}}^{\circ}, \Phi_{+, \omega^{\prime}}^{\circ}$ are taken to be disjoint, respectively, considering slightly larger convex cones. We check $\left(\Theta_{\omega^{\prime}}^{\circ}, \Theta_{\omega}\right)$-cone hyperbolicity of $F_{\alpha, \omega \omega^{\prime}}$, recalling Definition 2.2.2. The supports of $\widetilde{\Phi}_{-, \omega^{\prime}}^{\circ}, \widetilde{\Phi}_{+, \omega^{\prime}}, \widetilde{\Phi}_{0, \omega^{\prime}}^{\circ}$ are chosen analogously on corresponding slightly larger cones. The first compact inclusion in $(2.13)$ is a direct consequence of the compact inclusion in (2.17). To see the second compact inclusion in (2.13) note that

$$
\left(\left(\mathrm{D} F_{-\alpha, \omega \omega^{\prime}}\right)^{\operatorname{tr}}\right)^{-1}\left(\mathbb{R}^{d} \backslash \operatorname{supp} \Phi_{+, \omega}\right)=\mathbb{R}^{d} \backslash\left(\mathrm{D} F_{\alpha, \omega^{\prime} \omega}\right)^{\operatorname{tr}} \operatorname{supp} \Phi_{+, \omega}
$$

Comparing with the compact inclusion in (2.16), we conclude. The claim in (2.14) follows again by comparing with the compact inclusions in (2.16) and (2.17).

### 2.3 The transfer operator and the anisotropic Banach space

### 2.3.1 The transfer operator

We denote by $C^{r}(M)$ the space of $C^{\lfloor r\rfloor}$ functions whose $\lfloor r\rfloor$-th partial derivatives in charts are $C^{r-\lfloor r\rfloor}$. We let $C_{X}^{r-1}(M)^{3}$ be the space of $C^{r-1}$ functions which are $C^{r}$ in the flow direction $X$ defined by (2.4). Fixing a "potential function" $V \in C^{r-1}(M, \mathbb{R})$, we introduce the $\phi_{\alpha}$-weighted transfer operator family

$$
\begin{equation*}
\mathcal{L}_{\alpha, \phi_{\alpha}}: \varphi \mapsto \phi_{\alpha} \cdot\left(\varphi \circ g_{-\alpha}\right), \quad \alpha \geqslant 0 \tag{2.18}
\end{equation*}
$$

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acting on $\varphi \in C_{X}^{r-1}(M)$, where

$$
\phi_{\alpha}(x):=\exp \left(\int_{0}^{\alpha} V \circ g_{-\alpha^{\prime}}(x) \mathrm{d} \alpha^{\prime}\right) .
$$

We will construct Banach spaces $W_{p}^{s, t, q}$ containing $C_{X}^{r-1}(M)$ as a dense subspace (for suitable choices $p, s, t, q \in \mathbb{R}$ ) on which the family (2.18) of operators extends continuously to a strongly continuous semigroup (see Lemma 2.4.4 below). Note that

$$
\begin{equation*}
V=\partial_{\alpha} \phi_{\alpha \mid \alpha=0^{+}} . \tag{2.19}
\end{equation*}
$$

Our construction will show that for all $\varphi \in C_{X}^{r-1}(M)$

$$
\partial_{\alpha} \mathcal{L}_{\alpha, \phi_{\alpha}} \varphi_{\mid \alpha=0^{+}}=X \varphi+V \varphi,
$$

is well-defined in the sense that $(X+V) \varphi \in W_{p}^{s, t, q}$ if $\varphi \in C_{X}^{r-1}(M)$. The operator $X+V$ is the generator of the semigroup $\left\{\mathcal{L}_{\alpha, \phi_{\alpha}}: W_{p}^{s, t, q} \rightarrow W_{p}^{s, t, q} \mid \alpha \geqslant 0\right\}$. We denote by

$$
\left.\sigma(X+V)\right|_{W_{p}^{s, t, q}}
$$

the spectrum of $X+V$ to emphasize the dependency of the domain and hence the spectrum of $X+V$ on $W_{p}^{s, t, q}$. We show in Theorem 2.4.5 that the resolvent of $X+V$

$$
\begin{equation*}
\mathcal{R}_{z} \varphi:=(z-V-X)^{-1} \varphi,\left.\quad z \notin \sigma(X+V)\right|_{W_{p}^{s, t, q}}, \quad \varphi \in W_{p}^{s, t, q}, \tag{2.20}
\end{equation*}
$$

admits a Lasota-Yorke inequality for large $\Re z>0$. This allows us to identify a vertical left-open strip in the complex plane in which $\left.\sigma(X+V)\right|_{W_{p}^{s, t, q}}$ contains only isolated eigenvalues of finite multiplicity of $X+V$ (see Lemma 2.4.10).

### 2.3.2 The anisotropic Banach space

We work locally with the atlas $\mathcal{A}$, introduced in Section 2.2 . We let $\Psi_{n}: \mathbb{R}^{d} \rightarrow$ $[0,1], n \in \mathbb{Z}_{\geqslant 0}$, be a Paley-Littlewood decomposition as follows:
Let $\chi: \mathbb{R}_{>0} \rightarrow[0,1]$ be a $C^{\infty}$ map so that $\chi_{\mid(0,1]} \equiv 1$ and supp $\chi \subseteq[0,2]$. Let $|\cdot|: \mathbb{R}^{d} \rightarrow \mathbb{R}_{\geqslant 0}$ be a smooth norm on $\mathbb{R}^{d} \backslash\{0\}$. Define $\Psi_{n}$ by setting for all
$\xi \in \mathbb{R}^{d} \backslash\{0\}$

$$
\begin{equation*}
\Psi_{0}(\xi):=\chi(|\xi|) \quad \text { and } \quad \Psi_{n}(\xi):=\chi\left(\left|2^{-n} \xi\right|\right)-\chi\left(\left|2^{1-n} \xi\right|\right), n \geqslant 1 . \tag{2.21}
\end{equation*}
$$

This defines a partition of unity on $\mathbb{R}^{d} \backslash\{0\}$ since we have

$$
\sum_{n=0}^{\infty} \Psi_{n}(\xi)=\lim _{n \rightarrow \infty} \chi\left(\left|2^{-n} \xi\right|\right)=1
$$

For all $n \geqslant 1$ it holds $\Psi_{n}(\xi)=\Psi_{1}\left(2^{-n+1} \xi\right)$ from which one finds

$$
\begin{equation*}
\operatorname{supp} \Psi_{n} \subseteq\left\{\xi \in \mathbb{R}^{d}\left|2^{n-1} \leqslant|\xi| \leqslant 2^{n+1}\right\}\right. \tag{2.22}
\end{equation*}
$$

The inverse Fourier transform is given by

$$
\mathbb{F}^{-1} \varphi(x):=(2 \pi)^{-d} \int_{\mathbb{R}^{d}} e^{\mathrm{i} \xi x} \varphi(\xi) \mathrm{d} \xi
$$

where

$$
\xi x:=\langle\xi, x\rangle
$$

is the canonical scalar product on $\mathbb{R}^{d}$. The convolution of two complex valued functions $\varphi_{1}, \varphi_{2}$ on $\mathbb{R}^{d}$ (and extended to distributions) is given by

$$
\varphi_{1} * \varphi_{2}(x):=\int_{\mathbb{R}^{d}} \varphi_{1}(x-y) \varphi_{2}(y) \mathrm{d} y
$$

We will make frequent use (e.g. in the proof of Lemma 2.3.1 below and Lemma 2.4.16 in Section 2.4.4) of a special case of Young's inequality for convolutions,

$$
\left\|\varphi_{1} * \varphi_{2}\right\|_{L_{p}} \leqslant\left\|\varphi_{1}\right\|_{L_{1}}\left\|\varphi_{2}\right\|_{L_{p}}, \quad \text { for all } p \in[1, \infty] .
$$

Given a cone ensemble $\Theta=\left(\Phi_{-}, \Phi_{+}, \Phi_{0}\right)$, we set for all $\sigma \in\{-,+, 0\}, n \in \mathbb{Z}_{\geqslant 0}$,

$$
\begin{equation*}
\Psi_{\sigma, n}:=\Psi_{n} \Phi_{\sigma} \quad \text { and } \quad \Psi_{\sigma, n}^{\mathrm{Op}} \varphi:=\left(\mathbb{F}^{-1} \Psi_{\sigma, n}\right) * \varphi \tag{2.23}
\end{equation*}
$$

We let $\widetilde{\Psi}_{0}, \widetilde{\Psi}_{1} \in C^{\infty}$ such that $\widetilde{\Psi}_{0 \mid \operatorname{supp} \Psi_{0}} \equiv 1$ and $\widetilde{\Psi}_{1 \mid \operatorname{supp} \Psi_{1}} \equiv 1$. We set for every $n \in \mathbb{N}$

$$
\widetilde{\Psi}_{n}:=\widetilde{\Psi}_{1} \circ 2^{1-n} .
$$

(In principal it is enough to require the condition on the support of $\widetilde{\Psi}_{n}$ for each $n$ individually. Regarding the bounds in (2.25) below our choice here is

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reasonable.) Then we set for every $\sigma \in\{-,+, 0\}$ and every $n \in \mathbb{Z}_{\geqslant 0}$

$$
\begin{equation*}
\widetilde{\Psi}_{\sigma, n}:=\widetilde{\Psi}_{n} \widetilde{\Phi}_{\sigma} \quad \text { and } \quad \widetilde{\Psi}_{\sigma, n}^{\mathrm{Op}} \varphi:=\left(\mathbb{F}^{-1} \widetilde{\Psi}_{\sigma, n}\right) * \varphi, \tag{2.24}
\end{equation*}
$$

where $\widetilde{\Phi}_{\sigma} \in C^{\infty}$ and $\operatorname{supp} \tilde{\Phi}_{\sigma}$ is a closed convex cone such that $\tilde{\Phi}_{\sigma \mid \operatorname{supp} \Phi_{\sigma}} \equiv 1$ and $\Phi_{\sigma_{1}} \Phi_{\sigma_{2}} \equiv 0 \Rightarrow \widetilde{\Phi}_{\sigma_{1}} \widetilde{\Phi}_{\sigma_{2}} \equiv 0$ for all $\sigma_{1}, \sigma_{2} \in\{-,+, 0\}$. We have the following estimates for all $\sigma$ and all $n \in \mathbb{N}$ :

$$
\begin{equation*}
\left\|\mathbb{F}^{-1} \Psi_{n}\right\|_{L_{1}}=\left\|\mathbb{F}^{-1} \Psi_{1}\right\|_{L_{1}}<\infty, \quad\left\|\mathbb{F}^{-1} \Psi_{\sigma, n}\right\|_{L_{1}}=\left\|\mathbb{F}^{-1} \Psi_{\sigma, 1}\right\|_{L_{1}}<\infty . \tag{2.25}
\end{equation*}
$$

Analogous estimates hold for $\mathbb{F}^{-1} \Psi_{\sigma, 0}$ and $\mathbb{F}^{-1} \Psi_{0}$ and for the $\sim$-versions as well. If $\Theta^{\circ}$ is another cone ensemble we define $\Psi_{\sigma, n}^{\circ}, \Psi_{\sigma, n}^{\circ}$ Op and $\tilde{\Psi}_{\sigma, n}^{\circ}, \widetilde{\Psi}_{\sigma, n}^{\circ}{ }^{\circ \mathrm{p}}$ analogously. We set

$$
\begin{equation*}
B:=\left\{x \in \mathbb{R}^{d}| | x \mid<1\right\} \quad \text { and } \quad B^{\mathrm{c}}:=\mathbb{R}^{d} \backslash B . \tag{2.26}
\end{equation*}
$$

In order to show a continuous embedding of certain spaces we will use very often the following statement about convolution operators (an extension of [62, Theorem 0.3.1] for the case $r=1$ and $K(x, y)=K(x-y)$ in his notation). In Lemma 2.3.1 below all the occurring $L_{p}$-spaces are understood (as Bochner spaces, cf. [16]) such that if $a \in L_{p}\left(\mathbb{R}^{d}, \mathcal{B}\right)$ for some complex Banach space $\mathcal{B}$ then the norm of $a$ is given by

$$
\|a\|_{L_{p}\left(\mathbb{R}^{d}, \mathcal{B}\right)}:=\| \| a\left\|_{\mathcal{B}}\right\|_{L_{p}\left(\mathbb{R}^{d}, \mathbb{R}_{\geqslant 0}\right)} .
$$

The following lemma handles the range $p \in[1, \infty]$. (For parameters $p \in(1, \infty)$ one could apply instead the classical Marcinkiewicz theorem quoted e.g. as [10, Theorem 3.1].)

Lemma 2.3.1. Let $\mathcal{B}_{1}$ and $\mathcal{B}_{2}$ be (complex) Banach spaces, let $d \in \mathbb{N}$ and let $Q \in C^{d+1}\left(\mathbb{R}^{d}, \mathcal{L}\left(\mathcal{B}_{1}, \mathcal{B}_{2}\right)\right)$ satisfy for its partial derivatives

$$
\left\|\partial_{\xi}^{\beta} Q(\xi)\right\|_{\mathcal{L}\left(\mathcal{B}_{1}, \mathcal{B}_{2}\right)} \leqslant C(\beta)|\xi|^{-|\beta|} \quad \text { as } \quad|\xi| \rightarrow \infty
$$

for some constants $C(\beta)>0$ and all multi-indices $\beta \in\{0, \ldots, d+1\}^{d}$ such that $|\beta| \leqslant d+1$, where $|\beta|:=\beta_{1}+\ldots+\beta_{d}$. Then for all $p \in[1, \infty]$ the map

$$
\begin{equation*}
Q^{\mathrm{Op}}: L_{p}\left(\mathbb{R}^{d}, \mathcal{B}_{1}\right) \rightarrow L_{p}\left(\mathbb{R}^{d}, \mathcal{B}_{2}\right): a \mapsto \int_{\mathbb{R}^{d}}\left(\mathbb{F}^{-1} Q\right)(x-y) a(y) \mathrm{d} y \tag{2.27}
\end{equation*}
$$

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defines a bounded linear operator, where for every $b \in \mathcal{B}_{1}$ and every $x \in \mathbb{R}^{d}$

$$
\mathbb{F}^{-1} Q(x) b:=(2 \pi)^{-d} \int_{\mathbb{R}^{d}} e^{\mathrm{i} x \xi} Q(\xi) b \mathrm{~d} \xi
$$

It holds

$$
\left\|Q^{\mathrm{Op}}\right\|_{\mathcal{L}\left(L_{p}\left(\mathbb{R}^{d}, \mathcal{B}_{1}\right), L_{p}\left(\mathbb{R}^{d}, \mathcal{B}_{2}\right)\right)} \leqslant\left\|\mathbb{F}^{-1} Q\right\|_{L_{1}\left(\mathbb{R}^{d}, \mathcal{L}\left(\mathcal{B}_{1}, \mathcal{B}_{2}\right)\right)}<\infty
$$

Proof. Linearity of $Q^{\mathrm{Op}}$ follows if $Q^{\mathrm{Op}}$ is a bounded operator. Suppose first that $\mathbb{F}^{-1} Q \in L_{1}\left(\mathbb{R}^{d}, \mathcal{L}\left(\mathcal{B}_{1}, \mathcal{B}_{2}\right)\right)$. We estimate

$$
\begin{aligned}
\left\|Q^{\mathrm{Op}} a\right\|_{L_{p}\left(\mathbb{R}^{d}, \mathcal{B}_{2}\right)} & =\| \| \int_{\mathbb{R}^{d}}\left(\mathbb{F}^{-1} Q\right)(\cdot-y) a(y) \mathrm{d} y\left\|_{\mathcal{B}_{2}}\right\|_{L_{p}} \\
& \leqslant\left\|\int_{\mathbb{R}^{d}}\right\|\left(\mathbb{F}^{-1} Q\right)(\cdot-y) a(y)\left\|_{\mathcal{B}_{2}} \mathrm{~d} y\right\|_{L_{p}} \\
& \leqslant\left\|\int_{\mathbb{R}^{d}}\right\|\left(\mathbb{F}^{-1} Q\right)(\cdot-y)\left\|_{\mathcal{L}\left(\mathcal{B}_{1}, \mathcal{B}_{2}\right)}\right\| a(y)\left\|_{\mathcal{B}_{1}} \mathrm{~d} y\right\|_{L_{p}} \\
& =\| \| \mathbb{F}^{-1} Q\left\|_{\mathcal{L}\left(\mathcal{B}_{1}, \mathcal{B}_{2}\right)} *\right\| a\left\|_{\mathcal{B}_{1}}\right\|_{L_{p}} .
\end{aligned}
$$

Using Young's inequality, we estimate and conclude

$$
\begin{aligned}
\left\|\left\|\mathbb{F}^{-1} Q\right\|_{\mathcal{L}\left(\mathcal{B}_{1}, \mathcal{B}_{2}\right)} *\right\| a\left\|_{\mathcal{B}_{1}}\right\|_{L_{p}} & \leqslant\| \| \mathbb{F}^{-1} Q\left\|_{\mathcal{L}\left(\mathcal{B}_{1}, \mathcal{B}_{2}\right)}\right\|_{L_{1}}\| \| a\left\|_{\mathcal{B}_{1}}\right\|_{L_{p}} \\
& =\left\|\mathbb{F}^{-1} Q\right\|_{L_{1}\left(\mathbb{R}^{d}, \mathcal{L}\left(\mathcal{B}_{1}, \mathcal{B}_{2}\right)\right)}\|a\|_{L_{p}\left(\mathbb{R}^{d}, \mathcal{B}_{1}\right)} .
\end{aligned}
$$

We now show $\left\|\mathbb{F}^{-1} Q\right\|_{L_{1}\left(\mathbb{R}^{d}, \mathcal{L}\left(\mathcal{B}_{1}, \mathcal{B}_{2}\right)\right)}<\infty$. It remains to show an upper bound for

$$
I:=\int_{\mathbb{R}^{d}}\left\|\int_{\mathbb{R}^{d}} e^{\mathrm{i} y \xi} Q(\xi) a \mathrm{~d} \xi\right\|_{\mathcal{B}_{2}} \mathrm{~d} y
$$

Inside $I$ we substitute, whenever $y \neq 0$

$$
\xi \mapsto\langle y, y\rangle^{-\frac{1}{2}} \xi
$$

which yields

$$
I=\int_{\mathbb{R}^{d}}\langle y, y\rangle^{-\frac{d}{2}}\|J(y)\|_{\mathcal{B}_{2}} \mathrm{~d} y
$$

where

$$
J:=J(y)=\int_{\mathbb{R}^{d}} e^{\mathrm{i} y\langle y, y\rangle^{-\frac{1}{2}} \xi} Q\left(\langle y, y\rangle^{-\frac{1}{2}} \xi\right) a \mathrm{~d} \xi
$$

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For every $y \in \mathbb{R}^{d} \backslash\{0\}$ we set

$$
\xi_{0}:=y\langle y, y\rangle^{-\frac{1}{2}} \pi .
$$

Clearly, it holds

$$
\left\langle\xi_{0}, \xi_{0}\right\rangle=\pi^{2} .
$$

We now repeat the following substitution $(d+1)$-times

$$
\xi \mapsto \xi+\xi_{0},
$$

which yields

$$
J=\int_{\mathbb{R}^{d}} e^{\mathrm{i} y\langle y, y\rangle^{-\frac{1}{2}} \xi} \widetilde{Q}(\xi) a \mathrm{~d} \xi,
$$

where

$$
\widetilde{Q}(\xi)=2^{-d-1} \sum_{n=0}^{d+1}\binom{d+1}{n}(-1)^{n} Q\left(\langle y, y\rangle^{-\frac{1}{2}}\left(\xi+n \xi_{0}\right)\right) .
$$

We let $0<\epsilon<1$. We split the part in $\widetilde{Q}$ if $y \in B$ for every $0 \leqslant n \leqslant d+1$ as

$$
1 \equiv \chi\left(\left(\xi+n \xi_{0}\right)\langle y, y\rangle^{\frac{\epsilon}{2}}\right)+(1-\chi)\left(\left(\xi+n \xi_{0}\right)\langle y, y\rangle^{\frac{\epsilon}{2}}\right),
$$

for every corresponding summand in $\widetilde{Q}$, respectively. The part in $I$ which corresponds to $\chi\left(\left(\xi+n \xi_{0}\right)\langle y, y\rangle^{\frac{\epsilon}{2}}\right)$ is estimated trivially, using boundedness of $Q$ and integrability of $\langle y, y\rangle^{\frac{\epsilon-d}{2}}$ on $B$. Using the identity

$$
\begin{aligned}
Q(\xi)-Q\left(\xi+\xi_{0}\right) & =\int_{0}^{1} \partial_{t} Q\left(\left(\xi+\xi_{0}-t \xi_{0}\right)\langle y, y\rangle^{-\frac{\epsilon}{2}}\right) \mathrm{d} t \\
& =-\langle y, y\rangle^{-\frac{\epsilon}{2}} \int_{0}^{1}(\mathrm{D} Q)\left(\left(\xi+\xi_{0}-t \xi_{0}\right)\langle y, y\rangle^{-\frac{\epsilon}{2}}\right) \xi_{0} \mathrm{~d} t
\end{aligned}
$$

we now write the remaining part in $\widetilde{Q}$ as

$$
\begin{align*}
& \widetilde{Q}(\xi)=2^{-d-1} \int_{[0,1]^{d+1}} \sum_{n=0}^{d+1}\binom{d+1}{n}\langle y, y\rangle^{\frac{n \epsilon}{2}}\langle y, y\rangle^{-\frac{d+1-n}{2}}  \tag{2.28}\\
& \times\left(\left(\mathrm{D}^{n}(1-\chi)\right) \circ\langle y, y\rangle^{\frac{\epsilon}{2}}\left(\mathrm{D}^{d+1-n} Q\right) \circ\langle y, y\rangle^{-\frac{1}{2}}\right)(\xi(t))\left(-\xi_{0}\right)^{\otimes(d+1)} \mathrm{d} t,
\end{align*}
$$

where we put

$$
\xi(t):=\xi+(d+1) \xi_{0}-t \xi_{0}^{\otimes(d+1)} \quad \text { and } \quad \xi_{0}^{\otimes(d+1)}:=\underbrace{\left(\xi_{0}, \ldots, \xi_{0}\right)}_{(d+1)-\text { times }} .
$$

### 2.3 The transfer operator and the anisotropic Banach space

We observe that the part where derivatives of $\left((1-\chi) \circ\langle y, y\rangle^{\frac{\epsilon}{2}}\right)(\xi)$ contribute implies

$$
\xi \in((2 B) \backslash B)\langle y, y\rangle^{-\frac{\epsilon}{2}} .
$$

Using the decay condition on all the partial derivatives of $Q$, recalling that

$$
\langle y, y\rangle^{\frac{\epsilon-d}{2}} \log \langle y, y\rangle
$$

is integrable on the unit ball, and exchanging the order of integration with respect to $t$ as the outermost (justified by absolute integrability), we find for the corresponding part in $I$, for some constants $C_{1}, C_{2}, C_{3}>0$

$$
\begin{aligned}
I & \leqslant C_{1} \frac{\pi^{d+1}}{2^{d+1}}\|a\|_{\mathcal{B}_{1}} \sum_{n=0}^{d+1}\binom{d+1}{n} \int_{B} \int_{((2 B) \backslash B)\langle y, y\rangle^{-\frac{\epsilon}{2}}}|\xi|^{-d-1+n} \mathrm{~d} \xi\langle y, y\rangle^{\frac{n \epsilon-d}{2}} \mathrm{~d} y \\
& +C_{1} \pi^{d+1}\|a\|_{\mathcal{B}_{1}} \int_{B} \int_{2\langle y, y\rangle^{-\frac{\epsilon}{2}} B^{c}}|\xi|^{-d-1} \mathrm{~d} \xi\langle y, y\rangle^{-\frac{d}{2}} \mathrm{~d} y \\
& \leqslant C_{2} \pi^{d+1}(\log d)\|a\|_{\mathcal{B}_{1}} \int_{B}\left(1+\frac{\epsilon}{2}|\log \langle y, y\rangle|\right)\langle y, y\rangle^{\frac{\epsilon-d}{2}} \mathrm{~d} y \\
& \leqslant C_{3}\|a\|_{\mathcal{B}_{1}}
\end{aligned}
$$

In the case $y \in B^{\mathrm{c}}$ we proceed analogously, using the formular for $\widetilde{Q}(\xi)$ given in (2.28), but without splitting the integral with respect to $\xi$. We have now

$$
\widetilde{Q}(\xi)=2^{-d-1} \int_{[0,1]^{d+1}}\langle y, y\rangle^{-\frac{d+1}{2}}\left(\mathrm{D}^{d+1} Q\right)\left(\langle y, y\rangle^{-\frac{1}{2}} \xi(t)\right)\left(-\xi_{0}\right)^{\otimes(d+1)} \mathrm{d} t .
$$

If $\xi \in B\langle y, y\rangle^{\frac{\epsilon}{2}}$ we bound the corresponding part in $I$ trivially, using boundedness of the $(d+1)$-th partial derivatives of $Q$ and integrability of $\langle y, y\rangle_{\frac{d \epsilon-2 d-1}{2}}^{2}$ on $B^{\mathrm{c}}$. If $\xi \in B^{\mathrm{c}}\langle y, y\rangle^{\frac{\epsilon}{2}}$ we use the decay condition of the $(d+1)$-th partial derivatives of $Q$ instead and integrability of $\langle y, y\rangle^{-\frac{d+\epsilon}{2}}$ on $B^{\mathrm{c}}$.

For every open set $K \subseteq \mathbb{R}^{d}$ with compact closure we let $C_{0}^{r-1}(K)$ be the space of $C^{r-1}$ functions which vanish at the boundary of $K$. Since $C_{0}^{r-1}(K) \subset L_{p}(K, \mathbb{C})$ for all $p \in[1, \infty]$, the following definition makes sense.

Definition 2.3.2 (Local norm and local Banach space). Let $p \in[1, \infty]$ and let $s, t, q<r-1$. Let $\Theta$ be a cone ensemble from Definition 2.2.1 and let $K \subset \mathbb{R}^{d}$ be an open set with compact closure. For every $\varphi \in C_{0}^{r-1}(K)$ we set as the local

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norm

$$
\|\varphi\|_{W_{p, \Theta, K}^{s, t, q}}:=\left\|\left(\sum_{n=0}^{\infty} 4^{n s}\left|\Psi_{-, n}^{\mathrm{Op}} \varphi\right|^{2}+4^{n t}\left|\Psi_{+, n}^{\mathrm{Op}} \varphi\right|^{2}+4^{n q}\left|\Psi_{0, n}^{\mathrm{Op}} \varphi\right|^{2}\right)^{\frac{1}{2}}\right\|_{L_{p}\left(\mathbb{R}^{d}\right)} .
$$

The completion $W_{p, \Theta, K}^{s, t, q}$ of $C_{0}^{r-1}(K)$ under $\|\cdot\|_{W_{p, \theta, K}^{s, t, q}}$ is our local anisotropic Banach space.

This is an anisotropic version of a Triebel-Lizorkin space [63, p.45, Definition 2] with a certain inner $l_{2}$-norm and an outer $L_{p}$-norm. More precisely, we relate the summation in $n$ and $\sigma$ which appears in the norm of $W_{p, \Theta, K}^{s, t, q}$ to the norm of a Hilbert space of complex valued sequences defined on $\{-,+, 0\} \times \mathbb{Z} \geqslant 0$. We set

$$
\begin{equation*}
c(-):=s, \quad c(+):=t, \quad c(0):=q . \tag{2.29}
\end{equation*}
$$

Then we denote by $\ell_{2}^{c}$ the Hilbert space with norm given for all $a \in \ell_{2}^{c}$ by

$$
\begin{equation*}
\|a\|_{\ell_{2}^{c}}:=\left(\sum_{\sigma, n} 4^{c(\sigma) n}\left|a_{\sigma, n}\right|^{2}\right)^{\frac{1}{2}} \tag{2.30}
\end{equation*}
$$

For $s^{\prime}, t^{\prime}, q^{\prime} \in \mathbb{R}$ we define $c^{\prime}$ and $\ell_{2}^{c}$ analogously.
Lemma 2.3.3 (Multiplication and composition operator). Let $p \in[1, \infty]$ and let $s^{\prime}, t^{\prime}, q^{\prime}, s, t, q<r-1$. Let $\tilde{r}>\max \{0, s, t, q\}-\min \left\{0, s^{\prime}, t^{\prime}, q^{\prime}\right\}$ and let $f \in C_{0}^{\tilde{r}}(K)$ for some open set $K \in \mathbb{R}^{d}$ with compact closure and let $F: K \rightarrow F(K)$ be a $C^{\tilde{r}}$ diffeomorphism. Let $\Theta$ and $\Theta^{\circ}$ be two cone ensembles. Then the linear operator

$$
\mathcal{M}_{F, f}: W_{p, \Theta^{\circ}, K}^{s^{\prime}, t^{\prime},,^{\prime}} \rightarrow W_{p, \Theta, F^{-1}(K)}^{s, t, q}: \varphi \mapsto f \cdot(\varphi \circ F)
$$

is bounded if $c(\sigma) \leqslant c^{\prime}(\tau)$ whenever $\bigcup_{x \in K} \operatorname{supp} \Psi_{\sigma} \cap \mathrm{D} F(x)^{\operatorname{tr}} \operatorname{supp} \Psi_{\tau}^{\circ} \neq \varnothing$. Moreover, if $F=\mathrm{id}$ and $\Theta=\Theta^{\circ}$ the linear operator $\mathcal{M}_{\mathrm{id}, f}$ is bounded if $s \leqslant q \leqslant$ $t$.

Proof. We exclude first the indices for given $\sigma, \tau \in\{-,+, 0\}$ such that

$$
\begin{equation*}
\bigcup_{x \in K} \operatorname{supp} \Psi_{\sigma} \cap \mathrm{D} F(x)^{\operatorname{tr}} \operatorname{supp} \Psi_{\tau}^{\circ} \neq \varnothing, \tag{2.31}
\end{equation*}
$$

and given $n, \ell \in \mathbb{Z}_{\geqslant 0}$ such that

$$
\left.\sup _{x \in F(K), \xi \in\left(\operatorname{supp} \Psi_{\sigma} \cap B\right)}\left|\mathrm{D} F^{-1}(x)^{\operatorname{tr}} \xi\right|\right|^{-1} 2^{-4} \leqslant 2^{n-l} \leqslant 2^{4} \sup _{x \in K, \eta \in\left(\operatorname{supp} \Psi_{\tau}^{o} \cap B\right)}\left|\mathrm{D} F(x)^{\operatorname{tr}} \eta\right| .
$$

For all remaining $\sigma, \tau \in\{-,+, 0\}$ and $n, \ell \in \mathbb{Z}_{\geqslant 0}$ we bound the local norm for every $\epsilon>0$ and some constant $C_{1}=C_{1}(\epsilon)>0$

$$
\begin{align*}
\left\|\mathcal{M}_{F, f} \varphi\right\|_{W_{p, \Theta, K}^{s, t, q}} & =\left\|\left(\sum_{\sigma, n} 4^{-\epsilon n} 4^{(c(\sigma)+\epsilon) n}\left|\Psi_{\sigma, n}^{\mathrm{Op}} \mathcal{M}_{F, f} \varphi\right|^{2}\right)^{\frac{1}{2}}\right\|_{L_{p}\left(\mathbb{R}^{d}\right)} \\
& \leqslant C_{1} \sup _{\sigma, n} 2^{(c(\sigma)+\epsilon) n}\left\|\Psi_{\sigma, n}^{\mathrm{Op}_{\mathrm{p}}} \mathcal{M}_{F, f} \varphi\right\|_{L_{p}\left(\mathbb{R}^{d}\right)} \tag{2.32}
\end{align*}
$$

On the excluded indices we estimate as in the proof of Lemma 2.4.13 below, using Lemma 2.3.1 and Cauchy-Schwarz in $\ell$ and that $n \sim \ell$ and using $c(\sigma) \leqslant c^{\prime}(\tau)$. We recall the map $\widetilde{\Psi}_{\tau, \ell}^{\circ}$ defined in (2.24). Then we bound for every $n \geqslant 0$ and every $\sigma \in\{-,+, 0\}$

$$
\begin{align*}
\left\|\Psi_{\sigma, n}^{\mathrm{Op}} \mathcal{M}_{F, f} \varphi\right\|_{L_{p}\left(\mathbb{R}^{d}\right)} & \leqslant \sum_{\tau, \ell}\left\|\Psi_{\sigma, n}^{\mathrm{Op}} \mathcal{M}_{F, f} \tilde{\Psi}_{\tau, \ell}^{\circ \mathrm{Op}} \Psi_{\tau, \ell}^{\circ \mathrm{Op}} \varphi\right\|_{L_{p}\left(\mathbb{R}^{d}\right)} \\
& =\sum_{\tau, \ell} 2^{-c^{\prime}(\tau) \ell} 2^{c^{c}(\tau) \ell} \| \Psi_{\sigma, n}^{\mathrm{Op}} \mathcal{M}_{F, f} \widetilde{\Psi}_{\tau, \ell}^{\circ \mathrm{Op}} \Psi_{\tau, \ell}^{\circ} \mathrm{Op} \tag{2.33}
\end{align*} \|_{L_{p}\left(\mathbb{R}^{d}\right)} .
$$

Note that if $\operatorname{supp} \Psi_{\sigma}$ and $\mathrm{D} F(x)^{\operatorname{tr}} \operatorname{supp} \Psi_{\tau}^{\circ}$ have empty intersection, since the supports are open, we may assume that supp $\Psi_{\sigma}$ and $\mathrm{D} F(x)^{\operatorname{tr}} \operatorname{supp} \tilde{\Psi}_{\tau}^{\circ}$ have empty intersection as well. Since we excluded the conditions regarding certain $\sigma, \tau$ and $n, \ell$ given in (2.31) and below of it then by construction of $\tilde{\Psi}_{\tau, \ell}^{\circ}$, for some constant $C_{2}>0$ it holds, in the following assuming $n, \ell>0$,

$$
\begin{gather*}
\inf _{x \in K}\left|\operatorname{supp} \Psi_{\sigma, n}-\mathrm{D} F(x)^{\operatorname{tr}} \operatorname{supp} \tilde{\Psi}_{\tau, \ell}^{\circ}\right| \geqslant C_{2} 2^{\max \{n, \ell\}} \text { or } \\
\inf _{x \in F(K)}\left|\mathrm{D} F^{-1}(x)^{\operatorname{tr}} \operatorname{supp} \Psi_{\sigma, n}-\operatorname{supp} \tilde{\Psi}_{\tau, \ell}^{\circ}\right| \geqslant C_{2} 2^{\max \{n, \ell\}} \tag{2.34}
\end{gather*}
$$

In the following we assume the first inequality in (2.34). Otherwise the next estimates are done with the substitution $F(y) \mapsto y$. If $n=0$ or $\ell=0$ the following estimate is done analogously, using that either $\xi$ or $\eta$ is bounded. We set

$$
\tilde{\xi}:=2^{-n} \xi, \quad \tilde{\eta}:=2^{-\ell} \eta, \quad \text { and } \quad U:=\mathbb{R}^{d} \times \mathbb{R}^{d} \times K \times \mathbb{R}^{d} .
$$

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We write for every $x \in \mathbb{R}$

$$
\begin{aligned}
I(x) & :=I_{\sigma, n, \tau, \ell}(x):=\frac{(2 \pi)^{2 d}}{2^{(n+\ell) d}} \Psi_{\sigma, n}^{\mathrm{Op}} \mathcal{M}_{F, f} \Psi_{\tau, \ell}^{\circ \mathrm{Op}} \varphi(x) \\
& =\frac{(2 \pi)^{2 d}}{2^{(n+\ell) d}} \Psi_{\sigma, n}^{\mathrm{Op}} \mathcal{M}_{F, f} \widetilde{\Psi}_{\tau, \ell}^{\circ \mathrm{Op}} \Psi_{\tau, \ell}^{\circ \mathrm{Op}} \varphi(x) \\
& =\int_{U} e^{\mathrm{i} 2^{n} \widetilde{\xi}(x-y)} e^{\mathrm{i} 2^{\ell} \widetilde{\eta}(F(y)-z)} \Psi_{\sigma, 1}(\widetilde{\xi}) \widetilde{\Psi}_{\tau, 1}^{\circ}(\widetilde{\eta}) f(y) \Psi_{\tau, \ell}^{\circ \mathrm{Op}} \varphi(z) \mathrm{d} z \mathrm{~d} y \mathrm{~d} \tilde{\xi} \mathrm{~d} \widetilde{\eta} .
\end{aligned}
$$

Note that by assumption we have

$$
\begin{equation*}
\tilde{r}>\max \{0, s, t, q\}-\min \left\{0, s^{\prime}, t^{\prime}, q^{\prime}\right\} \geqslant 0 . \tag{2.35}
\end{equation*}
$$

Integrating $\tilde{r}$-times by parts (see Lemma B.3-Lemma B.5) in $y$, using the lower bound in (2.33), we arrive at
$I(x)=\int_{U} e^{\mathrm{i} 2^{n} \widetilde{\xi}(x-y)} e^{\mathrm{i} 2^{\ell} \tilde{\eta}(F(y)-z)} \Psi_{\sigma, 1}(\widetilde{\xi}) \widetilde{\Psi}_{\tau, 1}^{\circ}(\widetilde{\eta}) \frac{f_{\widetilde{r}}(y)}{2^{\max \{n, \ell\} \widetilde{r}}} \Psi_{\tau, \ell}^{\circ \mathrm{Op}} \varphi(z) \mathrm{d} z \mathrm{~d} y \mathrm{~d} \tilde{\xi} \mathrm{~d} \widetilde{\eta}$,
where all derivatives of $f_{\widetilde{r}}(y)$ with respect to $\tilde{\eta}$ and $\tilde{\xi}$ are bounded uniformly for all $(\tilde{\xi}, \tilde{\eta}, y) \in \operatorname{supp} \Psi_{\sigma, 1} \times \operatorname{supp} \tilde{\Psi}_{\tau, 1}^{\circ} \times K$. We set for every $y \in \mathbb{R}^{d}$ and for every $n \geqslant 0$

$$
u(y):=\left\{\begin{array}{ll}
1, & |y| \leqslant 1 \\
|y|^{-d-1}, & \text { otherwise }
\end{array}, \quad u_{n}:=u \circ 2^{n}\right.
$$

If $|x-y| 2^{n}>1$ we integrate $(d+1)$-times by parts in $\tilde{\xi}$ and if $|z-F(y)| 2^{\ell}>1$ we integrate $(d+1)$-times by parts in $\tilde{\eta}$. Hence we arrive at
$I(x)=2^{-\max \{n, \ell\} \tilde{r}} \int_{U} \tilde{f}_{\widetilde{r}}(\tilde{\xi}, \tilde{\eta}, y) u_{n}(x-y) u_{\ell}(z-F(y)) \Psi_{\tau, \ell}^{\circ \mathrm{Op}} \varphi(z) \mathrm{d} z \mathrm{~d} y \mathrm{~d} \tilde{\xi} \mathrm{~d} \tilde{\eta}$,
where $\tilde{f}_{\widetilde{r}}(\tilde{\xi}, \tilde{\eta}, y)$ is uniformly bounded for all $(\tilde{\xi}, \tilde{\eta}, y) \in \operatorname{supp} \Psi_{\sigma, 1} \times \operatorname{supp} \tilde{\Psi}_{\tau, 1}^{\circ} \times$ $K$. Hence we estimate for some constant $C_{3}>0$

$$
\begin{equation*}
|I(x)| \leqslant C_{3} 2^{-\max \{n, \ell\} \tilde{r}} u_{n} *\left(u_{\ell} \circ F\right) *\left|\Psi_{\tau, \ell}^{\circ \mathrm{Op}} \varphi\right|(x) \tag{2.36}
\end{equation*}
$$

We estimate for every $\sigma, \tau \in\{-,+, 0\}$ and every $n, \ell \geqslant 1$, using the equality in (2.35) and assuming $\epsilon>0$ small enough,

$$
\begin{equation*}
2^{(c(\sigma)+\epsilon) n-c^{\prime}(\tau) \ell-\max \{n, \ell\} \widetilde{r}} \leqslant 2^{(\max \{s, t, q\}+\epsilon) n-\min \left\{s^{\prime}, t^{\prime}, q^{\prime}\right\} \ell-\max \{n, \ell\} \widetilde{r}} \leqslant 2^{-\epsilon \ell} \tag{2.37}
\end{equation*}
$$

Hence we bound, using the estimates in (2.32), (2.33), (2.36), two times Young's

### 2.3 The transfer operator and the anisotropic Banach space

inequality and the bound in (2.37), for some constants $C_{4}, \ldots, C_{6}>0$

$$
\begin{gathered}
\left\|\mathcal{M}_{F, f} \varphi\right\|_{W_{p, \theta, K}^{s, t, q}} \leqslant C_{1} \sup _{\sigma, n} \sum_{\tau, \ell} 2^{(c(\sigma)+\epsilon) n-c^{\prime}(\tau) \ell} 2^{c^{\prime}(\tau) \ell}\left\|\Psi_{\sigma, n}^{\mathrm{Op}} \mathcal{M}_{F, f} \widetilde{\Psi}_{\tau, \ell}^{\circ \mathrm{Op}} \Psi_{\tau, \ell}^{\circ \mathrm{Op}} \varphi\right\|_{L_{p}} \\
=C_{1}(2 \pi)^{-2 d} 2^{d(n+\ell)} \sup _{\sigma, n} \sum_{\tau, \ell} 2^{(c(\sigma)+\epsilon) n-c^{\prime}(\tau) \ell} 2^{c^{\prime}(\tau) \ell}\left\|I_{\sigma, n, \tau, \ell}\right\|_{L_{p}} \\
\leqslant C_{4} \sup _{\sigma, n} \sum_{\tau, \ell} 2^{(c(\sigma)+\epsilon) n-c^{\prime}(\tau) \ell-\max \{n, \ell\} \tilde{r}^{(n+\ell) d} 2^{c^{\prime}(\tau) \ell}\left\|u_{n} *\left(u_{\ell} \circ F\right) * \Psi_{\tau, \ell}^{\circ \mathrm{Op}} \varphi\right\|_{L_{p}}} \begin{array}{l}
\leqslant C_{5} \sum_{\tau, \ell} 2^{-\ell \epsilon} 2^{c^{\prime}(\tau) \ell}\left\|\Psi_{\tau, \ell}^{\circ \mathrm{Op}} \varphi\right\|_{L_{p}} \leqslant C_{6} \sup _{\tau, \ell} 2^{c^{\prime}(\tau) \ell}\left\|\Psi_{\tau, \ell}^{\circ \mathrm{Op}} \varphi\right\|_{L_{p}} .
\end{array} .
\end{gathered}
$$

To see the statement if $F=\mathrm{id}$ we estimate the corresponding cases $c^{\prime}(\tau)<c(\sigma)$ if $\sigma \neq \tau$ and $n \sim \ell$ and $n, \ell \neq 0$ in a different way. We use
$\widetilde{I}(x):=\int_{U} e^{\mathrm{i} 2^{n} \tilde{\xi}(x-y)} e^{\mathrm{i} 2^{\ell} \tilde{\eta}(y-z)} \widetilde{\Psi}_{\sigma, 1}(\widetilde{\xi}) \Psi_{\sigma, 1}(\widetilde{\xi}) \tilde{\Psi}_{\tau, 1}^{\circ}(\widetilde{\eta}) f(y) \Psi_{\tau, \ell}^{\mathrm{Op}} \varphi(z) \mathrm{d} z \mathrm{~d} y \mathrm{~d} \tilde{\xi} \mathrm{~d} \tilde{\eta}$.
We express $\Psi_{\sigma, 1}(\tilde{\xi})$, using the identity

$$
\Psi_{\sigma, 1}(\tilde{\xi})-\Psi_{\sigma, 1}(\tilde{\eta})=\int_{0}^{1}\left(\mathrm{D} \Psi_{\sigma, 1}\right)(\tilde{\xi}+(1-h)(\tilde{\eta}-\tilde{\xi})) \mathrm{d} h(\tilde{\xi}-\tilde{\eta}) .
$$

We repeat this $k$-times in the right-hand side of this identity, replacing $\tilde{\xi}$ and yielding in total $k+1$ terms. The first $k$ terms are linear combinations of

$$
\Psi_{\sigma, 1}(j \tilde{\eta}-(j-1) \tilde{\xi})
$$

where $1 \leqslant j \leqslant k+1$. If $j=0$ then this is just $\Psi_{\sigma, 1}(\tilde{\eta})$. The corresponding part in $\widetilde{I}(x)$ is hence

$$
\tilde{I}_{1}(x)=\widetilde{\Psi}_{\sigma, n}^{\mathrm{Op}}\left(f \cdot \Psi_{\sigma, n}^{\mathrm{Op}} \Psi_{\tau, \ell}^{\mathrm{Op}} \varphi\right)=\widetilde{\Psi}_{\sigma, n}^{\mathrm{Op}}\left(f \cdot \Psi_{\sigma, n}^{\mathrm{Op}}\left(1-\sum_{\sigma^{\prime} \neq \tau} \Psi_{\sigma^{\prime}, \ell}^{\mathrm{Op}}\right) \varphi\right) .
$$

Note that $\widetilde{\Psi}_{\sigma, n}$ and $\Psi_{\sigma, n}$ satisfy the vanishing conditions in Lemma 2.3.1 as seen as an operator $\ell_{2}^{c} \mapsto \ell_{2}^{c}$. Then we bound with some constant $C_{5}=C_{5}(f)$

$$
\left\|\tilde{I}_{1}\right\|_{L_{p}} \leqslant C_{5}\left\|\left(\sum_{n=0}^{\infty}\left|4^{\sigma n} \Psi_{\sigma, n}^{\mathrm{Op}} \varphi\right|^{2}\right)^{\frac{1}{2}}\right\|_{L_{p}}
$$

using two times Lemma 2.3.1 and that $c(-) \leqslant c(0) \leqslant c(+)$ and that $\Psi_{+,,}^{\mathrm{Op}_{,}} \Psi_{-, n}^{\mathrm{Op}} \equiv$

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0 . The terms where $j>0$ are dealt with, using first the substitution

$$
j \tilde{\eta}-(j-1) \tilde{\xi} \mapsto \tilde{\eta},
$$

and then $\widetilde{r}$-times integration by parts analogous as before. The last $k+1$ term is

$$
\widetilde{\Psi}_{k}(\tilde{\eta}):=\int_{[0,1]^{k}}\left(\mathrm{D}^{k} \Psi_{\sigma, 1}\right)\left(\tilde{\xi}+\sum_{j=1}^{k}\left(1-t_{j}\right)(\tilde{\eta}-\tilde{\xi})\right) \mathrm{d} t(\tilde{\xi}-\tilde{\eta})^{\otimes k}
$$

We split now according to the size $|\tilde{\eta}-\tilde{\xi}|$. We let $\epsilon>0$. We note that

$$
2^{-n \epsilon((d+1)-k)} \chi\left(|\tilde{\eta}-\tilde{\xi}| 2^{n \epsilon}\right) \widetilde{\Psi}_{k}(\tilde{\eta})
$$

satisfies the vanishing conditions in Lemma 2.3.1 uniformly in $\tilde{\xi}$ as seen as an operator $\ell_{2}^{c} \mapsto \ell_{2}^{c}$ in $\widetilde{\eta}$. We bound the $L_{p}$ norm of the corresponding part analogous as in the case $\tilde{I}_{1}$. This is bounded appropriately with the choice of $k$ below. On the range $\left(1-\chi\left(|\tilde{\eta}-\tilde{\xi}| 2^{n \epsilon}\right)\right)>0$ we integrate $\widetilde{r}$-times by parts in $y$ and then $(d+1)$-times in $\tilde{\xi}$ and $\tilde{\eta}$ in the corresponding part of $\widetilde{I}(x)$. The terms which depend on $\chi$ are treated as in the range $\chi\left(|\tilde{\eta}-\tilde{\xi}| 2^{n \epsilon}\right)>0$. In the remaining part we gained a factor $\sim 2^{(-n+\epsilon) \tilde{r}+n \epsilon(d+1)}$. We choose $\epsilon$ small compatible with the inequality given in (2.35) and then $k$ large enough such that

$$
2^{(c(\sigma)-c(\tau)) n} \leqslant 2^{-n \epsilon((2 d+2)-k)} .
$$

Lemma 2.3.4 (Continuity and compactness). Let $p \in[1, \infty]$, let $s^{\prime} \leqslant s, q^{\prime} \leqslant q$, $t^{\prime} \leqslant t$, and $s \leqslant q \leqslant t$ and let $\Theta, \Theta^{\circ}$ be two cone ensembles, recalling Definition 2.2.1. Suppose the compact inclusions

$$
\begin{equation*}
\operatorname{supp} \Phi_{0}^{\circ} \Subset \operatorname{supp} \Phi_{0} \cup \operatorname{supp} \Phi_{+} \quad \text { and } \quad \operatorname{supp} \Phi_{+}^{\circ} \Subset \operatorname{supp} \Phi_{+} . \tag{2.38}
\end{equation*}
$$

Then the inclusion

$$
W_{p, \Theta, K}^{s, t, q} \subseteq W_{p, \Theta^{\circ}, K}^{s^{\prime}, t^{\prime}, q^{\prime}}
$$

is continuous for every open subset $K \subset \mathbb{R}^{d}$ with compact closure. Moreover, if $s^{\prime}<s, t^{\prime}<t$ and $q^{\prime}<q$ then the inclusion

$$
W_{p, \Theta, K}^{s, t, q} \subseteq W_{p, \Theta, K}^{s^{\prime}, t^{\prime}, q^{\prime}}
$$

is compact.
Proof. We prove first the claim on the continuous inclusion. We set for all $n \in \mathbb{Z}_{\geqslant 0}$

$$
\begin{aligned}
F_{-, n} & :=2^{\left(s^{\prime}-s\right) n}\left(\Psi_{-, n}^{\circ} /\left(2^{(q-s) n} \Psi_{0, n}+2^{(t-s) n} \Psi_{+, n}+\Psi_{-, n}\right)\right), \\
F_{+, n} & :=2^{\left(t^{\prime}-t\right) n}\left(\Psi_{+, n}^{\circ} / \Psi_{+, n}\right) \\
F_{0, n} & :=2^{\left(q^{\prime}-q\right) n}\left(\Psi_{0, n}^{\circ} /\left(2^{(t-q) n} \Psi_{+, n}+\Psi_{0, n}\right)\right)
\end{aligned}
$$

We define a map $Q$ on the Hilbert space $\ell_{2}^{c}$ (with norm as given in (2.30)) by setting for all $\sigma \in\{-,+, 0\}, n \in \mathbb{Z}_{\geqslant 0}$ and all $a \in \ell_{2}^{c}$

$$
(Q a)_{\sigma, n}:=F_{\sigma, n} a_{\sigma, n}
$$

In Lemma 2.3.1 we take $\mathcal{B}_{1}=\mathcal{B}_{2}=\ell_{2}^{c}$. It follows from the definition of $\Psi_{\sigma, n}$ in (2.23), the compact inclusion assumptions in (2.38) and the assumptions on $s, t, q, s^{\prime}, t^{\prime}, q^{\prime}$ that $Q$ satisfies the decay conditions on $Q$ in Lemma 2.3.1. It follows that the corresponding operator $Q^{\mathrm{Op}}$ in (2.27) is bounded. Let $\varphi \in$ $W_{p, \Theta, K}^{s, t, q}$. We set for all $n \in \mathbb{N}_{\geqslant 0}$

$$
\begin{aligned}
b_{-, n} & :=2^{s n} \mathbb{F}^{-1}\left(2^{(q-s) n} \Psi_{0, n}+2^{(t-s) n} \Psi_{+, n}+\Psi_{-, n}\right) * \varphi, \\
b_{+, n} & :=2^{t n} \mathbb{F}^{-1} \Psi_{+, n} * \varphi \\
b_{0, n} & :=2^{q n} \mathbb{F}^{-1}\left(2^{(t-q) n} \Psi_{+, n}+\Psi_{0, n}\right) * \varphi
\end{aligned}
$$

Then $\left(b_{\sigma, n} \mid \sigma \in\{-,+, 0\}, n \in \mathbb{N} \geqslant 0\right)=: b \in L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{c}\right)$ by assumption on $\varphi$ and in particular it holds, for some constant $C \geqslant 1,\|b\|_{L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{c}\right)} \leqslant C\|\varphi\|_{W_{p, Q, K}^{s, t, q}}$. We estimate, using Lemma 2.3.1, and conclude

$$
\|\varphi\|_{W_{p, \theta^{\circ}, K}^{s^{\prime}, t^{\prime}, q^{\prime}}}=\| \| Q^{\mathrm{Op}} b\left\|_{\ell_{2}^{c}}\right\|_{L_{p}} \leqslant C\|b\|_{L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{c}\right)} .
$$

We show the claim on the compact inclusion. We let $U \subset W_{p, \Theta, K}^{s, t, q}$ be a bounded set in $W_{p, \Theta, K}^{s, t, q}$ with bound $R>0$. We set $c^{\prime}$ with respect to $s^{\prime}, t^{\prime}, q^{\prime}$ analogous to c. It is enough to find for each $\epsilon>0$ an open cover of $U$ in $W_{p, \Theta, K}^{s^{\prime}, t^{\prime}, q^{\prime}}$ where each open set in the cover has size $\sim \epsilon$. (This yields total boundedness of $U$ in $W_{p, \Theta, K}^{s^{\prime}, t^{\prime}, q^{\prime}}$ and hence compactness.) Now there is $\delta>0$ such that for all $\sigma \in\{-,+, 0\}$

$$
\begin{equation*}
c^{\prime}(\sigma)+\delta-c(\sigma)<0 . \tag{2.39}
\end{equation*}
$$

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For all $\varphi \in U$ and all $N \in \mathbb{N}$ we bound

$$
\begin{align*}
\left\|\sqrt{\sum_{\sigma, n \geqslant N} 4^{-\delta n}\left|2^{\left(c^{\prime}(\sigma)+\delta\right) n} \Psi_{\sigma, n}^{\mathrm{Op}} \varphi\right|^{2}}\right\|_{L_{p}} & \leqslant C \sup _{\sigma, n \geqslant N} 2^{\left(c^{\prime}(\sigma)+\delta\right) n}\left\|\Psi_{\sigma, n}^{\mathrm{Op}} \varphi\right\|_{L_{p}} \\
& \leqslant C 2^{\left(c^{\prime}(\sigma)+\delta-c(\sigma)\right) N} R, \tag{2.40}
\end{align*}
$$

for some $\sigma \in\{-,+, 0\}$. Recalling the bound in (2.39), we make the bound in (2.40) smaller than $\epsilon$ by taking $N=N(\epsilon, R)$ large enough. Suppose now that the embedding is not compact. Then there are infinitely many $\varphi_{m} \in U, m \in \mathbb{N}$, such that for all $m_{1}>m_{2}$ it holds

Recalling the bound in (2.40), it holds for some $n<N$ and some $\sigma \in\{-,+, 0\}$

$$
\begin{equation*}
\left\|\varphi_{m_{1}}-\varphi_{m_{2}}\right\|_{W_{p, \Theta_{\omega}, \kappa_{\omega}\left(V_{\omega}\right)}^{s^{\prime}, t^{\prime}, q^{\prime}}} \leqslant C 2^{\left(c^{\prime}(\sigma)+\delta\right) n}\left\|\Psi_{\sigma, n}^{\mathrm{Op}}\left(\varphi_{m_{1}}-\varphi_{m_{2}}\right)\right\|_{L_{p}} . \tag{2.42}
\end{equation*}
$$

Since $C_{0}^{r-1}(K)$ is dense in $W_{p, \Theta_{\omega}, \kappa_{\omega}, V_{\omega}}^{s^{\prime} t^{\prime}, V^{\prime}}$ we may assume $\varphi_{m} \in C_{0}^{r-1}(K)$. We set $S:=\cup_{\sigma, n<N} \operatorname{supp} \Psi_{\sigma, n}$. Since all $\varphi_{m}$ are uniformly bounded in $W_{p, \Theta_{\omega}, \kappa_{\omega}\left(V_{\omega}\right)^{s}}^{s, t, q}$ norm and $\operatorname{supp} \varphi_{m}$ is uniformly bounded in $m$ as well, the Fourier transform of $\varphi_{m}$ cannot diverge on a dense subset of $S$ as $m \rightarrow \infty$ (this would violate the Paley-Wiener Theorem [39, Theorem 1.7.7]). By passing to a subsequence in $m$ we may split $S=S_{1} \sqcup S_{2}$ such that the family $\left\{\mathbb{F} \varphi_{m \mid S_{1}} \mid m \in \mathbb{N}\right\}$ is uniformly bounded. Then, using again that $\varphi_{m}$ has compact support with maximal diameter independent of $m$, the family $\left\{\mathbb{F} \varphi_{m \mid S_{1}} \mid m \in \mathbb{N}\right\}$ is also uniformly equicontinuous. Hence by the Arzelà-Ascoli Theorem there is a subsequence in $m$ such that $\varphi_{m \mid S_{1}}$ is a Cauchy sequence in $C^{0}$. Repeating the argument inductively for the part $\mathbb{F} \varphi_{m \mid S_{2}}$, then using a diagonal argument, we find a subsequence in $m$ such that $\mathbb{F} \varphi_{m \mid S}$ is a Cauchy sequence in $C^{0}$. Hence the right-hand side in (2.42) can be made arbitrary small which contradicts the lower bound in (2.41) and we conclude.

Lemma 2.3.5 (Local derivative). Let $p \in[1, \infty], s, t, q<r-1$ and let $\varphi \in$ $W_{p, \Theta_{\omega}, \kappa_{\omega}\left(V_{\omega}\right)}^{s, t, q}$. It holds for some constant $C>0$, for every $1 \leqslant j \leqslant d$, for every $\sigma \in\{-,+, 0\}$ such that

$$
\xi_{j} \neq 0 \quad \text { if } \quad\left(\xi_{1}, \ldots, \xi_{d}\right)=\xi \in \operatorname{supp} \Psi_{1, \sigma}
$$

and for every $\tilde{r} \in \mathbb{R}$

$$
\begin{equation*}
\left\|\left(\sum_{n=0}^{\infty} 4^{\widetilde{r} n}\left|\Psi_{\sigma, n}^{\mathrm{Op}} \varphi\right|^{2}\right)^{\frac{1}{2}}\right\|_{L_{p}} \leqslant C\left\|\left(\sum_{n=0}^{\infty} 4^{(\widetilde{r}-1) n}\left|\Psi_{\sigma, n}^{\mathrm{Op}} \partial_{x_{j}} \varphi\right|^{2}\right)^{\frac{1}{2}}\right\|_{L_{p}} \tag{2.43}
\end{equation*}
$$

Proof. Using the triangle inequality, it is enough to consider only the terms with $n>0$. For every $\xi \in \mathbb{R}^{d} \backslash\{0\}$ and $b \in \mathbb{C}$ we put

$$
(\mathcal{D}(\xi) b)_{n}:=\mathrm{i} \frac{\xi_{j}}{2^{n}} \Psi_{\sigma, n}(\xi) b, \quad n \in \mathbb{N}
$$

We note

$$
\Psi_{\sigma, n}^{\mathrm{Op}}\left(\partial_{x_{j}} \varphi\right)=\left(\mathbb{F}^{-1} \Psi_{\sigma, n}\right) * \partial_{x_{j}} \varphi=\left(\partial_{x_{j}} \mathbb{F}^{-1} \Psi_{\sigma, n}\right) * \varphi=2^{n}\left(\mathcal{D}^{\mathrm{Op}} \varphi\right)_{n}
$$

We let $\ell_{2}$ be the space of complex valued sequences space over $\mathbb{N}$. As norm we set $\|a\|_{\ell_{2}}:=\sqrt{\sum_{n=1}^{\infty} 4^{\tilde{r} n}\left|a_{n}\right|^{2}}$. For every $\xi \in \mathbb{R}^{d} \backslash\{0\}$, every $a \in \ell_{2}$ and every $n \in \mathbb{N}$ we put

$$
(Q(\xi) a)_{n}:=-\mathrm{i} \frac{2^{n}}{\xi_{j}} \widetilde{\Psi}_{\sigma, n}(\xi) a_{n}
$$

where $\widetilde{\Psi}_{\sigma, n}$ is defined in (2.24). Note that $\left(Q^{\mathrm{Op}} \mathcal{D}^{\mathrm{Op}} \varphi\right)_{n}=\Psi_{\sigma, n}^{\mathrm{Op}} \varphi$. Moreover since $\xi_{j} \neq 0$ by assumption, the map $Q$ satisfies the decay condition on its derivatives as required in Lemma 2.3.1. Hence, using Lemma 2.3.1 with $\mathcal{B}_{1}=$ $\mathcal{B}_{2}=\ell_{2}$, the map $Q^{\mathrm{Op}}: L_{p}\left(\mathbb{R}^{d}, \ell_{2}\right) \rightarrow L_{p}\left(\mathbb{R}^{d}, \ell_{2}\right)$ is a bounded linear operator. We conclude, using the estimate for some constant $C>0$

$$
\left\|\left\|Q^{\mathrm{Op}} \mathcal{D}^{\mathrm{Op}} \varphi\right\|_{\ell_{2}}\right\|_{L_{p}} \leqslant C\| \| \mathcal{D}^{\mathrm{Op}} \varphi\left\|_{\ell_{2}}\right\|_{L_{p}}
$$

We recall the open cover $V_{\omega} \subseteq M$ and the chart maps $\kappa_{\omega} \in \mathcal{A}, \omega \in \Omega$, introduced in Section 2.2. Also we recall the vector space $C_{X}^{r-1}(M)$ from the beginning of Section 2.3.2.

Definition 2.3.6 (Anisotropic Banach space). Let $\vartheta_{\omega}: V_{\omega} \rightarrow[0,1]$ be a $C^{r}$ partition of unity adapted to the chart maps $\kappa_{\omega}$ and let $\Theta_{\omega}$ be hyperbolic cone ensembles, recalling Definition 2.2.2, where $\omega \in \Omega$. Let $p \in[1, \infty]$, let $s, q, t<$ $r-1$ and let $\alpha_{0}>0$. We put for every $\varphi \in C_{X}^{r-1}(M)$ and every $p \in[1, \infty]$

$$
\begin{equation*}
\|\varphi\|_{W_{p}^{s, t, q}}:=\left(\sum_{\omega \in \Omega} \int_{0}^{\alpha_{0}}\left\|\left(\vartheta_{\omega} \cdot\left(\mathcal{L}_{\alpha, \phi_{\alpha}} \varphi\right) \circ \kappa_{\omega}^{-1}\right)\right\|_{W_{p, \theta_{\omega}, \kappa_{\omega}\left(V_{\omega}\right)}^{2}}^{2} \mathrm{~d} \alpha\right)^{\frac{1}{2}} \tag{2.44}
\end{equation*}
$$

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We denote by $W_{p}^{s, t, q}$ the completion of $C_{X}^{r-1}(M)$ under this norm.
Remark 2.3.7. Note that $W_{p}^{s, t, q}$ depends on the dynamics, $\alpha_{0}$, the atlas $\mathcal{A}$ and the cone ensembles $\Theta_{\omega}, \omega \in \Omega$. We understand each $\vartheta_{\omega} \circ \kappa_{\omega}^{-1}$ in (2.44) as extended to $\mathbb{R}^{d}$ by zero. By Lemma 2.3.3 a $C^{r}$ change of the atlas and hence a change of the cone ensemble yields an equivalent norm if $s \leqslant q \leqslant t<r-1$. The integration with respect to $\alpha$ is a way to "project out" the small times where the flow is not sufficiently hyperbolic. This is similar to [27, Definition 8.1] and also Baladi-Liverani [9, p.705, (3.2)] with the supremum replaced by an integral in the latter case. In turn, for $p=2$ the space $W_{2}^{s, t, q}$ is a Hilbert space because the parallelogram law

$$
\left\|\varphi_{1}+\varphi_{2}\right\|_{W_{p}^{s, t, q}}^{2}+\left\|\varphi_{1}-\varphi_{2}\right\|_{W_{p}^{s, t, q}}^{2}=2\left\|\varphi_{1}\right\|_{W_{p}^{s, t, q}}^{2}+2\left\|\varphi_{2}\right\|_{W_{p}^{s, t, q}}^{2}
$$

holds [14, Proposition 15.2].
The compact inclusion of the local Banach space in Lemma 2.3.4 carries over to the anisotropic Banach space $W_{p}^{s, t, q}$.

Lemma 2.3.8 (Compactness). Let $p \in[1, \infty]$, let $s \leqslant q \leqslant t$ such that max $\{0, t\}-$ $\min \{0, s\}<r-1$ and let $s^{\prime}<s, t^{\prime}<t$ and $q^{\prime}<q$ such that $\max \left\{0, \min \left\{s^{\prime}, t^{\prime}, q^{\prime}\right\}\right\}-$ $\min \left\{0, s^{\prime}, t^{\prime}, q^{\prime}\right\}<r-1$. Then there exist cone ensembles $\Theta_{\omega}, \omega \in \Omega$, such that the inclusion

$$
W_{p}^{s, t, q} \subseteq W_{p}^{s^{\prime}, t^{\prime}, q^{\prime}}
$$

is compact.
Proof. We let $s^{\prime}<s, t^{\prime}<t, q^{\prime}<q$. Let $U \subseteq W_{p}^{s, t, q}$ be a bounded set in the norm of $W_{p}^{s, t, q}$. In order to show the compact inclusion we proceed analogous to the proof in Lemma 2.3.4. To this end we let $\varphi_{m} \in U, m \in \mathbb{N}$, be a sequence, satisfying the analog bound in (2.41). Suppose now that there has to be some fixed $\omega \in \Omega$ and some fixed $\alpha \geqslant 0$ such that there exists $C_{1}>0$ such that for all $m>0$

$$
\sum_{\omega \in \Omega}\left\|\left(\vartheta_{\omega} \cdot\left(\mathcal{L}_{\alpha, \phi_{\alpha}} \varphi_{m} \circ \kappa_{\omega}^{-1}\right)\right)\right\|_{W_{p, \Theta_{\omega}, \kappa_{\omega}\left(V_{\omega}\right)}^{s, t, q}} \leqslant C_{1} .
$$

and that there exists some $\epsilon>0$ such that for all $m_{1}>m_{2}$ (up to some subsequence)

$$
\begin{equation*}
\left\|\vartheta_{\omega} \cdot\left(\mathcal{L}_{\alpha, \phi_{\alpha}}\left(\varphi_{m_{1}}-\varphi_{m_{2}}\right)\right) \circ \kappa_{\omega}^{-1}\right\|_{W_{p, \Theta \omega, \kappa_{\omega}\left(V_{\omega}\right)}^{s^{\prime}, t^{\prime}, q^{\prime}}}>\epsilon . \tag{2.45}
\end{equation*}
$$

Since $\left(\vartheta_{\omega} \cdot\left(\mathcal{L}_{\alpha, \phi_{\alpha}} \varphi_{m}\right) \circ \kappa_{\omega}^{-1}\right) \in W_{p, \Theta_{\omega}, \kappa_{\omega}\left(V_{\omega}\right)}^{s^{\prime} t^{\prime}, q^{\prime}}$ we find a Cauchy subsequence,
using the statement on the compact inclusion in Lemma 2.3.4. Note that by the Mean Value Theorem there exist non-fixed $\alpha=\alpha(m)$ and $\alpha=\alpha\left(m_{1}, m_{2}\right)$ which satisfy these inequalities. In particular, we wish to find a Cauchy subsequence for the left-hand side in the inequality (2.45) for the choice $\alpha=\alpha\left(m_{1}, m_{2}\right)$. Suppose $0 \leqslant \alpha^{\prime} \leqslant \alpha_{0}$. We have

$$
\begin{aligned}
& \left\|\vartheta_{\omega} \cdot\left(\mathcal{L}_{2 \alpha_{0}, \phi_{2 \alpha_{0}}} \varphi_{m}\right) \circ \kappa_{\omega}^{-1}\right\|_{W_{p, \Theta_{,}, \kappa_{\omega}\left(V_{\omega}\right)}^{s, t, q}} \leqslant \\
& \sum_{\omega^{\prime} \in \Omega}\left\|\vartheta_{\omega} \cdot\left(\mathcal{L}_{2 \alpha_{0}-\alpha^{\prime}, \phi_{2 \alpha_{0}-\alpha^{\prime}}}\left(\vartheta_{\omega^{\prime}} \cdot \mathcal{L}_{\alpha^{\prime}, \phi_{\alpha^{\prime}}} \varphi_{m}\right)\right) \circ \kappa_{\omega}^{-1}\right\|_{W_{p, \Theta_{\omega}, \kappa_{\omega}\left(V_{\omega}\right)}^{s, t, q}} .
\end{aligned}
$$

By Lemma 2.2.3 there exists cone ensembles $\left(\Theta_{\omega}, \Theta_{\omega}^{\circ}\right), \omega \in \Omega$, satisfying the condition (2.38) in Lemma 2.3.4 such that the local diffeomorphism of $g_{-\alpha}$ with $\alpha \geqslant \alpha_{0}$ is cone hyperbolic. Then, using Lemma 2.3.3 and $s \leqslant q \leqslant t$, $\max \{0, t\}-$ $\min \{0, s\}<r-1$ and taking $\alpha^{\prime}=\alpha(m)$, we bound this sequence in $m$ uniformly from above. Let $s^{\prime \prime}=\min \left\{s^{\prime}, t^{\prime}, q^{\prime}\right\}$. Then, using Lemma 2.3.3, recalling that it holds $\max \left\{0, \min \left\{s^{\prime}, t^{\prime}, q^{\prime}\right\}\right\}-\min \left\{0, s^{\prime}, t^{\prime}, q^{\prime}\right\}<r-1$, we find (abusing the notation $\mathcal{L}_{\alpha, \phi_{\alpha}}$ with negative $\alpha$ )

$$
\begin{align*}
& \left\|\vartheta_{\omega} \cdot\left(\mathcal{L}_{\alpha^{\prime}, \phi_{\alpha^{\prime}}}\left(\varphi_{m_{1}}-\varphi_{m_{2}}\right)\right) \circ \kappa_{\omega}^{-1}\right\|_{W_{p, \Theta_{\omega}, \kappa_{\omega}\left(V_{\omega}\right)}^{s^{\prime \prime}, s^{\prime \prime}, s^{\prime \prime}}} \leqslant  \tag{2.46}\\
& \sum_{\omega^{\prime} \in \Omega}\left\|\left(\vartheta_{\omega} \cdot\left(\mathcal{L}_{\alpha^{\prime}-2 \alpha_{0}, \phi_{\alpha^{\prime}-2 \alpha_{0}}}\left(\vartheta_{\omega^{\prime}} \cdot \mathcal{L}_{2 \alpha_{0}, \phi_{2 \alpha_{0}}}\left(\varphi_{m_{1}}-\varphi_{m_{2}}\right)\right)\right) \circ \kappa_{\omega}^{-1}\right)\right\|_{W_{p, \Theta_{\omega}, \xi^{\prime \prime}\left(V_{\omega}\right)}^{s^{\prime \prime}, V^{\prime \prime}}} \\
& \leqslant C_{2} \sum_{\omega \in \Omega}\left\|\vartheta_{\omega} \cdot\left(\mathcal{L}_{2 \alpha_{0}, \phi_{2 \alpha_{0}}}\left(\varphi_{m_{1}}-\varphi_{m_{2}}\right)\right) \circ \kappa_{\omega}^{-1}\right\|_{W_{p, \theta_{\omega}, k_{\omega}\left(V_{\omega}\right)}^{s^{\prime}, t^{\prime}, q^{\prime}}} .
\end{align*}
$$

for some constant $C_{2}>0$ independent of the choice of $\alpha^{\prime}$ and of $m_{1}, m_{2}$. Now we take $\alpha^{\prime}=\alpha\left(m_{1}, m_{2}\right)$ and let the right-hand side vanish in (2.46) as $\left(m_{1}, m_{2}\right) \rightarrow$ $\infty$. Then from the left-hand side for all $\sigma, n$

$$
\left\|\Psi_{\sigma, n}^{\mathrm{Op}}\left(\vartheta_{\omega} \cdot\left(\mathcal{L}_{\alpha\left(m_{1}, m_{2}\right), \phi_{\alpha\left(m_{1}, m_{2}\right)}}\left(\varphi_{m_{1}}-\varphi_{m_{2}}\right)\right) \circ \kappa_{\omega}^{-1}\right)\right\|_{L_{p}} \rightarrow 0
$$

By uniform boundedness in $W_{p, \Theta_{\omega}, \kappa_{\omega}\left(V_{\omega}\right)}^{s, t, q}$ (analogous to the proof of Lemma 2.3.4), the lower bound in (2.45) comes from a finite number of terms

$$
\left\|\Psi_{\sigma, n}^{\mathrm{Op}}\left(\vartheta_{\omega} \cdot\left(\mathcal{L}_{\alpha\left(m_{1}, m_{2}\right), \phi_{\alpha\left(m_{1}, m_{2}\right)}}\left(\varphi_{m_{1}}-\varphi_{m_{2}}\right)\right) \circ \kappa_{\omega}^{-1}\right)\right\|_{L_{p}}
$$

with $n \leqslant C_{3}=C_{3}(\epsilon)$. Hence we found a Cauchy subsequence for the left-hand side in the inequality (2.45) for the choice $\alpha=\alpha\left(m_{1}, m_{2}\right)$.

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### 2.4 Properties of the transfer operator, the generator and its resolvent

### 2.4.1 Bounds on the transfer operator

We introduce a local transfer operator in (2.47) below and state a local norm estimate for this operator in Lemma 2.4.1. We then give a norm estimate for the transfer operator family (2.18) in Lemma 2.4.2, making use of Lemma 2.4.1. Let $K \subset \mathbb{R}^{d}$ be an open set. Let $f: \mathbb{R}^{d} \rightarrow \mathbb{C}$ be a $C_{0}^{r-1}(K)$-map and let $F: K \rightarrow F(K)$ be a $\left(\Theta^{\circ}, \Theta\right)$-cone hyperbolic $C^{r}$-diffeomorphism on $K$ (recall Definition 2.2.2). The $f$-weighted local transfer operator is defined by

$$
\begin{equation*}
L_{F, f}: C_{0}^{r-1}(F(K)) \rightarrow C_{0}^{r-1}(K): \varphi \mapsto f \cdot \varphi \circ F \tag{2.47}
\end{equation*}
$$

Recalling $\tilde{\Phi}_{-}^{\circ}, \tilde{\Phi}_{+}^{\circ}, \tilde{\Phi}_{0}^{\circ}$ from (2.24), we put for every subset $I \subseteq K$

$$
\begin{aligned}
& \|F\|_{-, I}:=\inf _{\substack{y \in I \\
0 \neq \eta \in \operatorname{supp} \tilde{\Phi}_{-}^{\circ}}} \frac{\left|\left(\mathrm{D}_{y} F\right)^{\operatorname{tr}} \eta\right|}{|\eta|},\|F\|_{+, I}:=\sup _{\substack{y \in I \\
0 \neq \eta \in \operatorname{supp} \widetilde{\Phi}_{+}^{\circ}}} \frac{\left|\left(\mathrm{D}_{y} F\right)^{\operatorname{tr}} \eta\right|}{|\eta|}, \\
& \|F\|_{0, I}:=\sup _{\substack{y \in I \\
0 \neq \eta \in \operatorname{supp} \\
\tilde{\Phi}_{0}^{\circ}}} \frac{\left|\left(\mathrm{D}_{y} F\right)^{\operatorname{tr}} \eta\right|}{|\eta|} .
\end{aligned}
$$

Lemma 2.4.1 (Upper bound for local transfer operator). Let $\{\mathcal{W}\}$ denote the connected components of $\operatorname{supp} f$. Let $p \in[1, \infty]$. Let

$$
s^{\prime}<s<0<q \leqslant t<r-1+s^{\prime}, \quad q^{\prime}<q, \quad t^{\prime}<t
$$

Then for every $\varphi \in W_{p, \Theta, F(K)}^{s, t, q}$ it holds

$$
\left\|L_{F, f} \varphi\right\|_{W_{p, \theta, K}^{s, t, q}} \leqslant C_{0}\|\varphi\|_{W_{p, \Theta^{\circ}, F(K)}^{s^{\prime}, t^{\prime}, q^{\prime}}}+C_{1}\|\varphi\|_{W_{p, \Theta^{\circ}, F(K)}^{s^{\prime}, t^{\prime}, q}}+C_{2}\|\varphi\|_{W_{p, \Theta^{\circ}, F(K)}^{s, t, q}}
$$

where, for some constants $C>0$ and $k>0$, it holds
$C_{0} \leqslant C \sum_{\mathcal{W}} \max \left\{1,\|F\|_{-, \mathcal{W}}^{1-r},\|F\|_{0, \mathcal{W}}^{1-r}\right\}\|\mathrm{D} F\|_{C^{r-1}(\mathcal{W})}^{k}\|f\|_{C^{r-1}(\mathcal{W})}\left\|\left.\operatorname{det} \mathrm{D} F\right|^{-\frac{1}{p}}\right\|_{L_{\infty}(\mathcal{W})}$,
$C_{1} \leqslant C \sup _{\mathcal{W}}\left\|f|\operatorname{det} \mathrm{D} F|^{-\frac{1}{p}}\right\|_{L_{\infty}(\mathcal{W})} \max \left\{1,\|F\|_{0, \mathcal{W}}^{q}\right\}$ and
$C_{2} \leqslant C \sup _{\mathcal{W}}\left\|f|\operatorname{det} \mathrm{D} F|^{-\frac{1}{p}}\right\|_{L_{\infty}(\mathcal{W})} \max \left\{\|F\|_{+, \mathcal{W}}^{t},\|F\|_{-, \mathcal{W}}^{s}\right\}$.
Lemma 2.4.1 is proven in Section 2.4.4. For every $s, t, \alpha \in \mathbb{R}$ and every $x \in M$
we set

$$
\begin{equation*}
\lambda^{(t, s, \alpha)}(x):=\max \left\{\left\|\left(\mathrm{D} g_{\alpha}\right)_{\mid E_{+, g_{-\alpha}(x)}^{*}}^{\operatorname{tr}}\right\|^{-t},\left\|\left(\mathrm{D} g_{-\alpha}\right)_{\mid E_{-, x}^{*}}^{\operatorname{tr}}\right\|^{s}\right\} . \tag{2.48}
\end{equation*}
$$

For $s<0<t$ this quantity decreases exponentially fast to 0 as $\alpha \rightarrow \infty$, which is a consequence of the Anosov property given in (2.3) of the flow $g_{\alpha}$.

Lemma 2.4.2 (Bound on the transfer operator). Let $p \in[1, \infty]$. Let

$$
s^{\prime}<s<0<q \leqslant t<r-1+s^{\prime} \quad \text { and } \quad t^{\prime}<t .
$$

There exist $\alpha_{0}>0$, cone ensembles $\Theta_{\omega}, \omega \in \Omega$, and constants $A>0$ and $C>0$, such that for all $\varphi \in W_{p}^{s, t, q}$ with $\|\varphi\|_{W_{p}^{s, t, q}}=1$ and all $\alpha \geqslant 0$ it holds

$$
\left\|\mathcal{L}_{\alpha, \phi_{\alpha}} \varphi\right\|_{W_{p}^{s, t, q}} \leqslant C e^{A \alpha}\|\varphi\|_{W_{p}^{s^{\prime}, t^{\prime}, q}}+C(\alpha+1)\left\|\phi_{\alpha}\left|\operatorname{det} \mathrm{D} g_{-\alpha}\right|^{-\frac{1}{p}} \cdot \lambda^{(t, s, \alpha)}\right\|_{L_{\infty}} .
$$

Proof. We recall the map $F_{-\alpha, \omega \omega^{\prime}}$ and the set $V_{\alpha, \omega \omega^{\prime}}$ defined in Lemma 2.2.3 for all $\alpha \geqslant 0$ and all $\omega^{\prime}, \omega \in \Omega$. By Lemma 2.2.3 there exist cone ensembles $\Theta_{\omega^{\prime}}^{\circ}, \Theta_{\omega}$ such that the map $F_{-\alpha, \omega \omega^{\prime}}$ is $\left(\Theta_{\omega^{\prime}}^{\circ}, \Theta_{\omega}\right)$-cone hyperbolic. We recall the partition of unity $\vartheta_{\omega}$ (see Definition 2.3.6). We let

$$
V_{\alpha, \omega \omega^{\prime}} \subseteq \tilde{V}_{\alpha, \omega \omega^{\prime}} \subseteq V_{\omega}
$$

such that $F_{-\alpha, \omega \omega^{\prime}}$ is also $\left(\Theta_{\omega^{\prime}}^{\circ}, \Theta_{\omega}\right)$-cone hyperbolic on $\kappa_{\omega}\left(\tilde{V}_{\alpha, \omega \omega^{\prime}}\right)$. This is possible due to the compact inclusion of cones as required in the cone-hyperbolicity definition. We let

$$
\vartheta_{\alpha, \omega \omega^{\prime}}: \tilde{V}_{\alpha, \omega \omega^{\prime}} \rightarrow[0,1]
$$

be a $C_{0}^{r-1}$ map such that

$$
\vartheta_{\alpha, \omega \omega^{\prime} \mid V_{\alpha, \omega \omega^{\prime}}} \equiv \vartheta_{\omega \mid V_{\alpha, \omega \omega^{\prime}}}
$$

For all $z \in \kappa_{\omega}\left(V_{\omega}\right)$ we have

$$
\vartheta_{\omega} \circ \kappa_{\omega}^{-1}(z) \cdot \vartheta_{\omega^{\prime}} \circ \kappa_{\omega^{\prime}}^{-1} \circ F_{-\alpha, \omega \omega^{\prime}}(z)=\vartheta_{\alpha, \omega \omega^{\prime}} \circ \kappa_{\omega}^{-1}(z) \cdot \vartheta_{\omega^{\prime}} \circ \kappa_{\omega^{\prime}}^{-1} \circ F_{-\alpha, \omega \omega^{\prime}}(z) .
$$

Note that $\left\|\vartheta_{\alpha, \omega \omega^{\prime}}\right\|_{C^{r-1}}$ is controlled by the rate of expansion of $F_{-\alpha, \omega \omega^{\prime}}$. Let $\varphi \in W_{p}^{s, t, q}$ and put $W_{\omega}:=W_{p, \Theta_{\omega}, \kappa_{\omega}\left(V_{\omega}\right)}^{s, t, q}$ and $W_{\omega \omega^{\prime}}:=W_{p, \Theta_{\omega}, \tilde{V}_{\alpha, \omega \omega^{\prime}}}^{s, t, q}$. For all

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$\alpha \geqslant \alpha_{0}$, for some $C \geqslant 1$, we estimate for every $p \in[1, \infty]$

$$
\begin{align*}
& \left\|\mathcal{L}_{\alpha, \phi_{\alpha}} \varphi\right\|_{W_{p}^{s, t, q}}^{2} \leqslant C \max _{\omega \in \Omega} \int_{0}^{\alpha_{0}}\left\|\left(\vartheta_{\omega} \cdot\left(\phi_{\alpha^{\prime}} \cdot \mathcal{L}_{\alpha, \phi_{\alpha}} \varphi \circ g_{-\alpha^{\prime}}\right)\right) \circ \kappa_{\omega}^{-1}\right\|_{W_{\omega}}^{2} \mathrm{~d} \alpha^{\prime} \\
& \quad=C \max _{\omega \in \Omega} \int_{0}^{\alpha}\left\|\left(\vartheta_{\omega} \cdot \phi_{\alpha}\right) \circ \kappa_{\omega}^{-1} \cdot \sum_{\omega^{\prime} \in \Omega}\left(\vartheta_{\omega^{\prime}} \cdot \mathcal{L}_{\alpha^{\prime}, \phi_{\alpha^{\prime}}} \varphi\right) \circ \kappa_{\omega^{\prime}}^{-1} \circ F_{-\alpha, \omega \omega^{\prime}}\right\|_{W_{\omega}}^{2} \mathrm{~d} \alpha^{\prime} \\
& \quad=C \max _{\omega \in \Omega} \int_{0}^{\alpha_{0}}\left\|\sum_{\omega^{\prime} \in \Omega}\left(\vartheta_{\alpha, \omega \omega^{\prime}} \cdot \phi_{\alpha}\right) \circ \kappa_{\omega}^{-1} \cdot\left(\vartheta_{\omega^{\prime}} \cdot \mathcal{L}_{\alpha^{\prime}, \phi_{\alpha^{\prime}}} \varphi\right) \circ \kappa_{\omega^{\prime}}^{-1} \circ F_{-\alpha, \omega \omega^{\prime}}\right\|_{W_{\omega}}^{2} \mathrm{~d} \alpha^{\prime} \\
& \leqslant C^{2} \max _{\omega, \omega^{\prime} \in \Omega} \int_{0}^{\alpha 0}\left\|\left(\vartheta_{\alpha, \omega \omega^{\prime}} \cdot \phi_{\alpha}\right) \circ \kappa_{\omega}^{-1} \cdot\left(\vartheta_{\omega^{\prime}} \cdot \mathcal{L}_{\alpha^{\prime}, \phi_{\alpha^{\prime}}} \varphi\right) \circ \kappa_{\omega^{\prime}}^{-1} \circ F_{-\alpha, \omega \omega^{\prime}}\right\|_{W_{\omega}}^{2} \mathrm{~d} \alpha^{\prime} \\
& =C^{2} \max _{\omega, \omega^{\prime} \in \Omega} \int_{0}^{\alpha_{0}}\left\|L_{F_{-\alpha, \omega \omega^{\prime}},\left(\vartheta_{\alpha, \omega \omega^{\prime}} \cdot \phi_{\alpha}\right) \circ \kappa_{\omega}^{-1}}\left(\left(\vartheta_{\omega^{\prime}} \cdot \mathcal{L}_{\alpha^{\prime}, \phi_{\alpha^{\prime}}} \varphi\right) \circ \kappa_{\omega^{\prime}}^{-1}\right)\right\|_{W_{\omega \omega^{\prime}}}^{2} \mathrm{~d} \alpha^{\prime} . \tag{2.49}
\end{align*}
$$

We used in the last step the definition of the weighted local transfer operator (see (2.47)) in which we take $F:=F_{-\alpha, \omega \omega^{\prime}}$ and as the $C_{0}^{r-1}$-weight $f:=\left(\vartheta_{\alpha, \omega \omega^{\prime}} \cdot \phi_{\alpha}\right) \circ$ $\kappa_{\omega}^{-1}$.
We now show the claimed upper bound for $\mathcal{L}_{\alpha, \phi_{\alpha}}$. We recall that

$$
\operatorname{supp} f=\kappa_{\omega}\left(\tilde{V}_{\alpha, \omega \omega^{\prime}}\right)=\bigsqcup \kappa_{\omega}(\mathcal{W}),
$$

where the disjoint union is over all the finitely many connected components $\mathcal{W}$ of $\tilde{V}_{\alpha, \omega \omega^{\prime}}$. The inclusion $W_{p, \Theta^{\circ}, \kappa_{\omega^{\prime}}\left(V_{\omega^{\prime}}\right)}^{s, t, q} \subseteq W_{\omega^{\prime}}$ is continuous by Lemma 2.3.4. Together with the bound given by Lemma 2.4.1 this yields the upper bound

$$
\left\|\mathcal{L}_{\alpha, \phi_{\alpha}} \varphi\right\|_{W_{p}^{s, t, q}} \leqslant \widetilde{C}_{1}\|\varphi\|_{W_{p}^{s^{\prime}, t^{\prime}, q}}+\widetilde{C}_{2}\|\varphi\|_{W_{p}^{s, t, q}},
$$

where
$\widetilde{C}_{1} \leqslant C^{2} \max _{\omega, \omega^{\prime} \in \Omega} C_{0}\left(F_{-\alpha, \omega \omega^{\prime}}, f\right)+C_{1}\left(F_{-\alpha, \omega \omega^{\prime}}, f\right), \widetilde{C}_{2} \leqslant C^{2} \max _{\omega, \omega^{\prime} \in \Omega} C_{2}\left(F_{-\alpha, \omega \omega^{\prime}}, f\right)$,
and $C_{0}, C_{1}, C_{2}$ are the constants from Lemma 2.4.1. We claim for some constant $C_{4}>0$ the following bound

$$
\begin{align*}
\left\|\phi_{\alpha}\left|\operatorname{det} \mathrm{D} g_{-\alpha}\right|^{-\frac{1}{p}}\right\|_{L_{\infty}(\mathcal{W})} & \leqslant\left. C_{4} \inf _{x \in \mathcal{W}}\left|\phi_{\alpha}\right| \operatorname{det} \mathrm{D} g_{-\alpha}\right|^{-\frac{1}{p}} \mid(x)  \tag{2.50}\\
& =C_{4}\left\|\left(\phi_{-\alpha}\left|\operatorname{det} \mathrm{D} g_{\alpha}\right|^{-\frac{1}{p}}\right) \circ g_{-\alpha}\right\|_{L_{\infty}(\mathcal{W})} .
\end{align*}
$$

Due to the construction of $\tilde{V}_{\alpha, \omega \omega^{\prime}}$, all points in a connected component $\mathcal{W}$ stay close under iterates by $g_{\alpha^{\prime}}$ for all $0 \leqslant \alpha^{\prime} \leqslant \alpha$. Then in the case of hyperbolic maps
the bound in (2.50) follows, using [43, Proposition 20.2.6.]. However for Anosov flows the distance between two points $x_{1}, x_{2} \in \mathcal{W}$ may never be sufficiently contracted under iterates by $g_{\alpha^{\prime}}$, e.g. if $x_{1}, x_{2}$ belong to a same orbit of $g_{\alpha^{\prime}}$. We split (along the flow direction $X$ in charts) each $\mathcal{W}$ into parts $\mathcal{W}_{j}, 1 \leqslant j \leqslant$ $\lfloor\alpha\rfloor+1$, in which now two points are no more than $\sim(\lfloor\alpha\rfloor+1)^{-1}$ apart. We set $\widetilde{\mathcal{W}}:=\left\{\mathcal{W}_{j}\right\}$ for all $1 \leqslant j \leqslant\lfloor\alpha\rfloor+1$. Then it holds the bound in (2.50) with $\mathcal{W}$ replaced by $\mathcal{W}_{j}$. We modify $\vartheta_{\alpha, \omega \omega^{\prime}}$, taking a sufficiently small neighborhood $U_{j}$ containing $\mathcal{W}_{j}$, such that $\vartheta_{\alpha, \omega \omega^{\prime} \mid U_{j}}$ is $C_{0}^{r}$. Then passing to this new weights $\vartheta_{\alpha, \omega \omega^{\prime} \mid U_{j}}$ and summing over $j$ we obtain an additional factor $\sim(\alpha+1)$ in the right-hand side in (2.49). We recall $\lambda^{(t, s, \alpha)}(x)$ from (2.48) and $\|F\|_{-, I},\|F\|_{+, I}$ introduced below (2.47) in which we take $I=\mathcal{W}_{j}$ and $F=F_{-\alpha, \omega \omega^{\prime}}$. In addition note $F_{-\alpha, \omega \omega^{\prime}}^{-1}=F_{\alpha, \omega^{\prime} \omega}$. Then we write
$\left\|F_{-\alpha, \omega \omega^{\prime}}\right\|_{+, \kappa_{\omega}\left(\mathcal{W}_{j}\right)}=\left(\inf _{\substack{y \in \eta \in\left(\mathrm{D}_{y} F_{-\alpha, \omega \omega^{\prime}}\left(\mathcal{W}_{j}\right) \\ \operatorname{tr} \\ \operatorname{supp} \tilde{\Phi}_{+, \omega^{\prime}}^{\circ}\right.}}^{|\eta|} \frac{\left|\left(\mathrm{D}_{F_{-\alpha, \omega \omega^{\prime}}(y)} F_{\alpha, \omega^{\prime} \omega}\right)^{\operatorname{tr}} \eta\right|}{|\eta|}\right)^{-1}$.
We recall the construction in (2.16) of the $C^{+}$-cones in the proof of Lemma 2.2.3. We find a compactly embedded cone

$$
C_{\gamma+, \omega}^{+} \Subset\left(\mathrm{D}_{y} F_{-\alpha, \omega \omega^{\prime}}\right)^{\operatorname{tr}} \operatorname{supp} \widetilde{\Phi}_{+, \omega^{\prime}}^{\circ}
$$

which is transversal to another cone $C_{\gamma-, \omega}^{-}$. Hence the unstable distribution $E_{-}$ (in charts) stays away from $\left(\mathrm{D}_{y} F_{-\alpha, \omega \omega^{\prime}}\right)^{\operatorname{tr}} \operatorname{supp} \tilde{\Phi}_{+, \omega^{\prime}}^{\circ}$ by some positive angle. Replacing the inf with the sup, it holds for some constant $C_{5}>0$

$$
\begin{equation*}
\left\|F_{-\alpha, \omega \omega^{\prime}}\right\|_{+, \kappa_{\omega}\left(\mathcal{W}_{j}\right)} \leqslant C_{5}\left(\sup _{x \in \mathcal{W}_{j}}\left\|\left(\mathrm{D} g_{\alpha}\right)_{\mid E_{+, g-\alpha}(x)}^{\mathrm{tr}}\right\|\right)^{-1} \tag{2.51}
\end{equation*}
$$

By analogous reasoning we conclude similar for $\left\|F_{-\alpha, \omega \omega^{\prime}}\right\|_{-, \kappa_{\omega}\left(\mathcal{W}_{j}\right)}$. We estimate

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for some constants $C_{6}, \ldots, C_{9}>0$, using the bounds in (2.51) and (2.50),

$$
\begin{aligned}
\widetilde{C}_{2} & \leqslant C_{6} \max _{\widetilde{\mathcal{W}}, \mathcal{W}_{j} \in \widetilde{\mathcal{W}}}\left\|\phi_{\alpha}\left|\operatorname{det} \mathrm{D} g_{-\alpha}\right|^{-\frac{1}{p}}\right\|_{L_{\infty}\left(\mathcal{W}_{j}\right)}\left\|\lambda^{(t, s, \alpha)}\right\|_{L_{\infty}\left(\mathcal{W}_{j}\right)} \\
& \leqslant \alpha C_{7} \max _{\widetilde{\mathcal{W}}, \mathcal{W}_{j} \in \widetilde{\mathcal{W}}}\left\|\phi_{\alpha}\left|\operatorname{det} \mathrm{D} g_{-\alpha}\right|^{-\frac{1}{p}}\right\|_{L_{\infty}\left(\mathcal{W}_{j}\right)}\left\|\lambda^{(t, s, \alpha)}\right\|_{L_{\infty}\left(\mathcal{W}_{j}\right)} \\
& \leqslant \alpha C_{8} \max _{\widetilde{\mathcal{W}}, \mathcal{W}_{j} \in \widetilde{\mathcal{W}}}\left\|\left(\phi_{-\alpha}\left|\operatorname{det} \mathrm{D} g_{\alpha}\right|^{-\frac{1}{p}}\right) \circ g_{-\alpha}\right\|_{L_{\infty}\left(\mathcal{W}_{j}\right)}^{-1}\left\|\lambda^{(t, s, \alpha)}\right\|_{L_{\infty}\left(\mathcal{W}_{j}\right)} \\
& \leqslant \alpha C_{9} \max _{\widetilde{\mathcal{W}}, \mathcal{W}_{j} \in \widetilde{\mathcal{W}}}\left\|\phi_{\alpha}\left|\operatorname{det} \mathrm{D} g_{-\alpha}\right|^{-\frac{1}{p}} \lambda^{(t, s, \alpha)}\right\|_{L_{\infty}\left(\mathcal{W}_{j}\right)} .
\end{aligned}
$$

Inspecting the constant $\widetilde{C}_{1}$, all terms depending on $F$ and $f$ are bounded by the maximal expansion of $F_{-\alpha, \omega \omega^{\prime}}$ and $\phi_{\alpha}$, respectively, which grow at most exponentially in $\alpha$. Hence, there is $A>0$ and $C_{10} \geqslant 1$ such that $\widetilde{C}_{1} \leqslant C_{11} e^{A \alpha}$. If $\alpha<\alpha_{0}$ we split $\int_{0}^{\alpha_{0}}=\int_{0}^{\alpha_{0}-\alpha}+\int_{\alpha_{0}-\alpha}^{\alpha_{0}}$. Hence it holds $\left\|\mathcal{L}_{\alpha, \phi_{\alpha}} \varphi\right\|_{W_{p}^{s, t, q}} \leqslant\|\varphi\|_{W_{p}^{s, t, q}}+$ $\left\|\mathcal{L}_{\alpha_{0}, \phi_{\alpha_{0}}} \varphi\right\|_{W_{p}^{s, t, q}}$. The latter term is estimated as in the case $\alpha \geqslant \alpha_{0}$. Since $\alpha \leqslant \alpha_{0}$, we combine here the upper bound of $\left\|\mathcal{L}_{\alpha, \phi_{\alpha}} \varphi\right\|_{W_{p}^{s, t, q}}$ with the second term of our desired estimate, increasing the constant $C_{11}$.

Remark 2.4.3. A weaker upper bound for the transfer operator, e.g.

$$
\left\|\mathcal{L}_{\alpha, \phi_{\alpha}}\right\|_{W_{p}^{s, t, q} \rightarrow W_{p}^{s, t, q}} \leqslant C_{1} \exp \left(C_{2} \alpha\right)
$$

for all $\alpha \geqslant 0$ and for some constants $C_{1}, C_{2} \geqslant 1$ independent of $\alpha$, can be obtained for a wider choice of $s, t, q$, e.g. for some $s>0$ (and this carries over to Lemma 2.4.4 below as well). However, we are interested in the parameter range as assumed in Lemma 2.4.2 which allows us to show the Lasota-Yorke inequality for the resolvent given in Theorem 2.4.5 below. See also Lemma 2.5.17 in the next section below for such a bound in the case of a special weight.
We recall that the family $\left\{\mathcal{L}_{\alpha, \phi_{\alpha}}: W_{p}^{s, t, q} \rightarrow W_{p}^{s, t, q} \mid \alpha \geqslant 0\right\}$ forms a strongly continuous semigroup if and only if $\lim _{\alpha \rightarrow 0^{+}}\left\|\mathcal{L}_{\alpha, \phi_{\alpha}} \varphi-\varphi\right\|_{W_{p}^{s, t, q}}=0$ for all $\varphi \in W_{p}^{s, t, q}$ (e.g. see [46, Proposition I.1.3]).

Lemma 2.4.4 (Strongly continuous semigroup). Let $p \in[1, \infty]$ and let $s<0<$ $q \leqslant t<r-1+s$. Then the transfer operator family

$$
\left\{\mathcal{L}_{\alpha, \phi_{\alpha}}: W_{p}^{s, t, q} \rightarrow W_{p}^{s, t, q} \mid \alpha \geqslant 0\right\}
$$

forms a strongly continuous semigroup.
2.4 Properties of the transfer operator, the generator and its resolvent

Proof. Let $\varphi \in W_{p}^{s, t, q}$. For fixed $s<0<q \leqslant t$ such that $t-s<r-1$ there is $\delta>0$ such that $t-s<r-1-\delta$. We set $s^{\prime}:=s-\delta$ and let $t^{\prime}<t$. Then $s, t, q, s^{\prime}$ and $t^{\prime}$ satisfy the assumptions of Lemma 2.4.2. Using Lemma 2.4.2, we bound the transfer operator for all small $\alpha \geqslant 0$

$$
\begin{equation*}
\left\|\mathcal{L}_{\alpha, \phi_{\alpha}} \varphi\right\|_{W_{p}^{s, t, q}} \leqslant C_{1}\|\varphi\|_{W_{p}^{s^{\prime}, t^{\prime}, q}}+C_{2}\|\varphi\|_{W_{p}^{s, t, q}} \leqslant\left(C_{1}+C_{2}\right)\|\varphi\|_{W_{p}^{s, t, q}} \tag{2.52}
\end{equation*}
$$

for some constants $C_{1}, C_{2}>0$ independent of $\alpha$. By density, for every $\epsilon>0$ there is $\tilde{\varphi} \in C_{X}^{r-1}(M)$ such that

$$
\begin{equation*}
\|\varphi-\widetilde{\varphi}\|_{W_{p}^{s, t, q}} \leqslant \epsilon \tag{2.53}
\end{equation*}
$$

Using first the triangle inequality and then the bounds (2.52)-(2.53), we estimate

$$
\begin{align*}
\left\|\mathcal{L}_{\alpha, \phi_{\alpha}} \varphi-\varphi\right\|_{W_{p}^{s, t, q}} & \leqslant\left\|\mathcal{L}_{\alpha, \phi_{\alpha}}(\varphi-\widetilde{\varphi})\right\|_{W_{p}^{s, t, q}}+\|\varphi-\widetilde{\varphi}\|_{W_{p}^{s, t, q}}+\left\|\mathcal{L}_{\alpha, \phi_{\alpha}} \widetilde{\varphi}-\widetilde{\varphi}\right\|_{W_{p}^{s, t, q}} \\
& \leqslant C_{3} \epsilon+\left\|\mathcal{L}_{\alpha, \phi_{\alpha}} \widetilde{\varphi}-\widetilde{\varphi}\right\|_{W_{p}^{s, t, q}} \tag{2.54}
\end{align*}
$$

for some constant $C_{3}>0$ independent of $\epsilon$ and $\alpha$. Since $\varphi \in C_{X}^{r-1}(M)$ we have

$$
\mathcal{L}_{\alpha, \phi_{\alpha}} \widetilde{\varphi}-\tilde{\varphi}=\alpha \int_{0}^{1}\left(\partial_{\alpha^{\prime}} \mathcal{L}_{\alpha^{\prime}, \phi_{\alpha^{\prime}}} \varphi\right)_{\mid \alpha^{\prime}=h \alpha} \mathrm{~d} h
$$

Since $\left(\partial_{\alpha^{\prime}} \mathcal{L}_{\alpha^{\prime}, \phi_{\alpha^{\prime}}} \varphi\right)_{\mid \alpha^{\prime}=h \alpha} \in C^{(r-1)}(M)$ the norm $\left\|\left(\partial_{\alpha^{\prime}} \mathcal{L}_{\alpha^{\prime}, \phi_{\alpha^{\prime}}} \varphi\right)_{\mid \alpha^{\prime}=h \alpha}\right\|_{W_{p}^{s, t, q}}$ is finite for all $0 \leqslant h \leqslant 1$. Hence for some constant $C_{4}(\varphi)=C_{4}>0$ we bound

$$
\begin{equation*}
\left\|\mathcal{L}_{\alpha, \phi_{\alpha}} \tilde{\varphi}-\widetilde{\varphi}\right\|_{W_{p}^{s, t, q}} \leqslant \alpha \sup _{0 \leqslant h \leqslant 1}\left\|\left(\partial_{\alpha^{\prime}} \mathcal{L}_{\alpha^{\prime}, \phi_{\alpha^{\prime}}} \varphi\right)_{\mid \alpha^{\prime}=h \alpha}\right\|_{W_{p}^{s, t, q}} \leqslant C_{4} \alpha \tag{2.55}
\end{equation*}
$$

We conclude by a combination of the estimates (2.54)-(2.55).

### 2.4.2 Lasota-Yorke inequality for the resolvent

We use Lemma 2.4.2 to prove Theorem 2.4.5 below. We use in addition that the resolvent improves regularity in the flow direction. We set, recalling $\lambda^{(t, s, \alpha)}$ in (2.48),

$$
\begin{equation*}
\lambda_{\min }=\lambda_{\min }(s, t, p):=\lim _{\alpha \rightarrow \infty} \frac{1}{\alpha} \log \left\|\phi_{\alpha}\left|\operatorname{det} \mathrm{D} g_{-\alpha}\right|^{-\frac{1}{p}} \lambda^{(t, s, \alpha)}\right\|_{L_{\infty}(M)} \tag{2.56}
\end{equation*}
$$

The following theorem will allow us to show that $\lambda_{\min }(s, t)$ plays the role of the essential spectral bound of $X+V$ :

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Theorem 2.4.5 (Lasota-Yorke inequality for the resolvent). Let $p \in[1, \infty]$ and let

$$
s^{\prime}<s<0<q \leqslant t<r-1+s^{\prime}, \quad q-1 \leqslant q^{\prime}<q, \quad t^{\prime}<t .
$$

There exist $\alpha_{0}>0, A_{0}>\lambda_{\text {min }}$, cone ensembles $\Theta_{\omega}, \omega \in \Omega$, and a constant $C>0$ such that for every $\varphi \in W_{p}^{s, t, q}$ with $\|\varphi\|_{W_{p}^{s, t, q}}=1$, for every $z \in \mathbb{C}$ with $\Re z>A_{0}$ and for every $n \in \mathbb{N}$ it holds

$$
\left\|\mathcal{R}_{z}^{n+1} \varphi\right\|_{W_{p}^{s, t, q}} \leqslant C \frac{|z|+1+\left(\Re z-A_{0}\right)}{\left(\Re z-A_{0}\right)^{n+1}}\|\varphi\|_{W_{p}^{s^{\prime}, t^{\prime}, q^{\prime}}}+\frac{C n\left(\Re z-\lambda_{\text {min }}\right)^{-1}+C}{\left(\Re z-A_{0}\right)\left(\Re z-\lambda_{\text {min }}\right)^{n}} .
$$

Proof. Since $\lambda^{(s, t, \alpha)}$ grows at most exponentially as $\alpha \rightarrow \infty$, the constant $\lambda_{\text {min }}$ is finite by a result on superadditive functions [37, Theorem 7.6.1]. We let $A_{0}>\lambda_{\text {min }}$. By Lemma 2.4.4 the transfer operator family (2.18) forms a strongly continuous semigroup with a well-defined generator $X+V$. We estimate powers of the resolvent $\mathcal{R}_{z}$ defined in (2.20). To this end we work with the integral representation of powers of the resolvent defined in (2.57) below (see [46, Corollary II.1.11]). We recall the constant $A$ given in Lemma 2.4.2 and let $A_{0}>A$. We set for every $z \in \mathbb{C}$ such that $\Re z>A_{0}$ and every $n \in \mathbb{N}$

$$
\begin{equation*}
\mathcal{R}_{z}^{n} \varphi:=\int_{0}^{\infty} \frac{\alpha^{n-1} e^{-z \alpha}}{(n-1)!} \mathcal{L}_{\alpha, \phi_{\alpha}} \varphi \mathrm{d} \alpha, \quad \varphi \in W_{p}^{s, t, q} . \tag{2.57}
\end{equation*}
$$

We have directly from (2.57) for all $\alpha \geqslant 0$

$$
\begin{equation*}
\mathcal{R}_{z}^{n} \mathcal{L}_{\alpha, \phi_{\alpha}} \varphi=\mathcal{L}_{\alpha, \phi_{\alpha}} \mathcal{R}_{z}^{n} \varphi . \tag{2.58}
\end{equation*}
$$

Using Lemma 2.4.2, we estimate for some constant $C_{1}>0$

$$
\begin{align*}
\left\|\mathcal{R}_{z}^{n+1} \varphi\right\|_{W_{p}^{s, t, q}} & \leqslant \int_{0}^{\infty} \frac{\alpha^{n-1} e^{-\Re z \alpha}}{(n-1)!}\left\|\mathcal{L}_{\alpha, \phi_{\alpha}} \mathcal{R}_{z} \varphi\right\|_{W_{p}^{s, t, q}} \mathrm{~d} \alpha \\
& \leqslant \frac{C_{1}}{\left(\Re z-A_{0}\right)^{n}}\left\|\mathcal{R}_{z} \varphi\right\|_{W_{p}^{s^{\prime}, t^{\prime}, q}}+\frac{C_{1}\left(n+\left(\Re z-\lambda_{\min }\right)\right)}{\left(\Re z-\lambda_{\min }\right)^{n+1}}\left\|\mathcal{R}_{z} \varphi\right\|_{W_{p}^{s, t, q}} . \tag{2.59}
\end{align*}
$$

Using Lemma 2.4.2, we get boundedness for some constant $C_{2}>0$

$$
\begin{equation*}
\left\|\mathcal{R}_{z} \varphi\right\|_{W_{p}^{s, t, q}} \leqslant \frac{C_{2}}{\Re z-A_{0}}\|\varphi\|_{W_{p}^{s, t, q}} . \tag{2.60}
\end{equation*}
$$

Therefore the second term in the right-hand side in (2.59) is bounded as claimed. We bound now the first term in the right-hand side in (2.59). Inverting the
flowbox condition (2.10), we find $\mathrm{D} \kappa_{\omega}^{-1} \partial_{x_{d}}=X_{\mid V_{\omega}}$. Hence it holds

$$
\begin{align*}
\partial_{x_{d}}\left(\vartheta_{\omega} \cdot \varphi\right) \circ \kappa_{\omega}^{-1}(x) & =\mathrm{D}\left(\vartheta_{\omega} \cdot \varphi\right) \circ \kappa_{\omega}^{-1}(x) \mathrm{D}_{x} \kappa_{\omega}^{-1} \partial_{x_{d}} \\
=\left(\mathrm{D}\left(\vartheta_{\omega} \cdot \varphi\right) X_{\mid V_{\omega}}\right) \circ \kappa_{\omega}^{-1}(x) & =\left(\left(X \vartheta_{\omega}\right) \cdot \varphi+\vartheta_{\omega} \cdot(X \varphi)\right) \circ \kappa_{\omega}^{-1}(x) . \tag{2.61}
\end{align*}
$$

We set $W_{\omega}^{q}:=W_{p, \Theta_{\omega}, \kappa_{\omega}\left(V_{\omega}\right)}^{s^{\prime}, t^{\prime},}, \omega \in \Omega$. We estimate the local norms inside the norm $\left\|\mathcal{R}_{z} \varphi\right\|_{W_{p}^{s^{\prime}, t^{\prime}, q},}$, using the equality in (2.61), then Lemma 2.3 .5 and the equality in (2.58), for some constant $C_{3}>0$ :

$$
\begin{array}{r}
\left\|\left(\vartheta_{\omega} \cdot \mathcal{L}_{\alpha^{\prime}, \phi_{\alpha^{\prime}}} \mathcal{R}_{z} \varphi\right) \circ \kappa_{\omega}^{-1}\right\|_{W_{\omega}^{q}} \leqslant C_{3}\left\|\left(\vartheta_{\omega} \cdot \mathcal{L}_{\alpha^{\prime}, \phi_{\alpha^{\prime}}} \varphi\right) \circ \kappa_{\omega}^{-1}\right\|_{W_{\omega}}^{q-1}  \tag{2.6}\\
+C_{3}\left\|\left(\left(X \vartheta_{\omega}\right) \cdot \mathcal{L}_{\alpha^{\prime}, \phi_{\alpha^{\prime}}} \mathcal{R}_{z} \varphi\right) \circ \kappa_{\omega}^{-1}\right\|_{W_{\omega}^{q-1}}+C_{3}\left\|\left(\vartheta_{\omega} \cdot X \mathcal{R}_{z} \mathcal{L}_{\alpha^{\prime}, \phi_{\alpha^{\prime}}} \varphi\right) \circ \kappa_{\omega}^{-1}\right\|_{W_{\omega}^{q-1}} .
\end{array}
$$

We note that $\left(X \vartheta_{\omega}\right) \circ \kappa_{\omega}^{-1} \in C_{0}^{r-1}\left(\kappa_{\omega}\left(V_{\omega}\right)\right)$ and $t-s<r-1$. Using Lemma 2.3.3, we bound for some constant $C_{4}(X)=C_{4}>0$

$$
\begin{equation*}
\left\|\left(\left(X v_{\omega}\right) \cdot \mathcal{L}_{\alpha^{\prime}, \phi_{\alpha^{\prime}}} \mathcal{R}_{z} \varphi\right) \circ \kappa_{\omega}^{-1}\right\|_{W_{\omega}^{q-1}}^{q-1} \leqslant C_{4} \sup _{\omega \in \Omega}\left\|\left(\vartheta_{\omega} \cdot \mathcal{L}_{\alpha^{\prime}, \phi_{\alpha^{\prime}}} \mathcal{R}_{z} \varphi\right) \circ \kappa_{\omega}^{-1}\right\|_{W_{\omega}^{q}}^{q-1} . \tag{2.63}
\end{equation*}
$$

Using the equality

$$
X \mathcal{R}_{z} \varphi=z \mathcal{R}_{z} \varphi-V \mathcal{R}_{z} \varphi-\varphi,
$$

together with the equality in (2.58), we find

$$
\begin{align*}
& \left\|\left(\vartheta_{\omega} \cdot X \mathcal{R}_{z} \mathcal{L}_{\alpha^{\prime}, \phi_{\alpha^{\prime}}} \varphi\right) \circ \kappa_{\omega}^{-1}\right\|_{W_{\omega}^{q-1}} \leqslant|z|\left\|\left(\vartheta_{\omega} \cdot \mathcal{L}_{\alpha^{\prime}, \phi_{\alpha^{\prime}}} \mathcal{R}_{z} \varphi\right) \circ \kappa_{\omega}^{-1}\right\|_{W_{\omega}^{q-1}}  \tag{2.64}\\
& \quad+\left\|\left(\vartheta_{\omega} \cdot V \mathcal{L}_{\alpha^{\prime}, \phi_{\alpha^{\prime}}} \mathcal{R}_{z} \varphi\right) \circ \kappa_{\omega}^{-1}\right\|_{W_{\omega}^{q-1}}+\left\|\left(\vartheta_{\omega} \cdot \mathcal{L}_{\alpha^{\prime}, \phi_{\alpha^{\prime}}} \varphi\right) \circ \kappa_{\omega}^{-1}\right\|_{W_{\omega}^{q-1}}^{q-1} .
\end{align*}
$$

Recalling that $V \in C^{r-1}(M)$, we bound the term which contains the factor $\left(\vartheta_{\omega} \cdot V\right) \circ \kappa_{\omega}^{-1}$ in the right-hand side in (2.64) analogous as in the estimate in (2.63). The final estimate follows by a combination of the bounds (2.59)(2.60) and (2.62)-(2.64), together with the trivial continuous inclusion $W_{\omega}^{q^{\prime}} \subseteq$ $W_{\omega}^{q-1}$.

A direct consequence of Theorem 2.4.5 is the bound on the essential spectral radius of the resolvent:

Corollary 2.4.6 (Essential spectral radius). Under the assumptions of Theorem 2.4 .5 (including the choices for $p, s, t, q \in \mathbb{R}$ ), letting $A_{0}$ and $\lambda_{\text {min }}=\lambda_{\min }(s, t, p)$
be the constants from that theorem, the essential spectral radius of the resolvent $\mathcal{R}_{z}: W_{p}^{s, t, q} \rightarrow W_{p}^{s, t, q}$ is bounded by $\left|\Re z-\lambda_{\text {min }}\right|^{-1}$ for all $z \in \mathbb{C}$ with $\Re z>A_{0}$.

Proof. Let $s^{\prime}<s, t^{\prime}<t$ and $q^{\prime}<q$. The inclusion $W_{p}^{s, t, q} \subseteq W_{p}^{s^{\prime}, t^{\prime}, q^{\prime}}$ is compact by Lemma 2.3.8. Then, together with a result of Hennion [36, Corollaire 1] and Theorem 2.4.5 we find the claimed bound on the essential spectral radius of the resolvent.

We recall $\lambda_{\text {min }}$ defined in (2.56) and $\phi_{\alpha}$ in (2.18).
Lemma 2.4.7. Let $d=3$ and let $\left|\operatorname{det} \mathrm{D} g_{\alpha}\right| \equiv 1$. Set $\tilde{t}:=\min \{-t, s\}$. Then it holds

$$
\lambda_{\text {min }}=\lim _{\alpha \rightarrow \infty} \frac{1}{\alpha} \log \left\|\phi_{\alpha}\left|\operatorname{det}\left(\mathrm{D} g_{-\alpha}\right)^{\operatorname{tr}}{ }_{\mid E_{-}^{*}}\right|\right\|_{L_{\infty}(M)} .
$$

Proof. Since the flow is volume preserving, we have

$$
\left.\left|\operatorname{det}\left(\mathrm{D} g_{-\alpha}\right)^{\operatorname{tr}}\right| E_{+}^{*}\right|^{-1}=\left|\operatorname{det}\left(\mathrm{D}_{g_{-\alpha}} g_{\alpha}\right)^{\operatorname{tr}}\right| E_{0}^{*}| | \operatorname{det}\left(\mathrm{D}_{g_{-\alpha}} g_{\alpha}\right)^{\operatorname{tr}}\left|E_{-}^{*}\right| .
$$

Since $d_{-}=1=d-2$ we can replace $\left\|\left(\mathrm{D} g_{-\alpha}\right)_{\left|E_{-, x}^{*}\right|}^{\operatorname{tr}}\right\|$ in $\lambda^{(s, t, \alpha)}$ by $\left|\operatorname{det}\left(\mathrm{D} g_{-\alpha}\right)^{\operatorname{tr}}{ }_{\mid E_{-, x}^{*}}\right|$ and $\left\|\left(\mathrm{D} g_{\alpha}\right)_{\mid E_{+, g_{-\alpha}(x)}^{*}}^{\operatorname{tr}}\right\|$ by $\left.\left|\operatorname{det}\left(\mathrm{D} g_{-\alpha}\right)^{\operatorname{tr}}\right| E_{+, x}^{*}\right|^{-1}$. Moreover $\left|\operatorname{det}\left(\mathrm{D}_{g_{-\alpha}} g_{\alpha}\right)^{\operatorname{tr}}\right| E_{0}^{*} \mid$ is bounded from above and below and we conclude.

Remark 2.4.8. Note that Lemma 2.4.7 holds in the particular case of a contact Anosov flow if $d=3$. Clearly, if $\left|\phi_{\alpha}\right| \leqslant\left.\left|\operatorname{det}\left(\mathrm{D} g_{-\alpha}\right)^{\operatorname{tr}}\right| E_{-}^{*}\right|^{-\widetilde{t}}$ for all $\alpha>0$ then $\lambda_{\text {min }} \leqslant 0<h_{\text {top }}$.

### 2.4.3 Spectral properties of the generator

All spectral properties of the generator $X+V$ are with respect to its domain $D(X+V)=D(X+V)_{\mid W_{p}^{s, t, q}}$ for admissible choices $p, s, t, q \in \mathbb{R}$ which is discussed in the following lemma.

Lemma 2.4.9 (Domain of the generator). Let $p, s, q, t \in \mathbb{R}$ satisfy the assumptions of Lemma 2.4.4. Then the family $\left\{\mathcal{L}_{\alpha, \phi_{\alpha}}: W_{p}^{s, t, q} \rightarrow W_{p}^{s, t, q} \mid \alpha \geqslant 0\right\}$ admits a generator

$$
X+V: D(X+V) \rightarrow W_{p}^{s, t, q}
$$

which is a closed operator on its domain $D(X+V)$. Moreover, the inclusion

$$
D(X+V) \subseteq W_{p}^{s, t, q}
$$

is dense and the inclusion

$$
C_{X}^{r-1}(M) \subseteq D(X+V)
$$

is dense for the graph norm $\|\cdot\|_{W_{p}^{s, t, q}}+\|(X+V)(\cdot)\|_{W_{p}^{s, t, q}}$.
Proof. Using Lemma 2.4.4, the statement about $X+V$ being a densely (in $W_{p}^{s, t, q}$ ) defined closed operator is [46, Theorem II.1.4]. Suppose now $\mathcal{L}_{\alpha, \phi_{\alpha}}\left(C_{X}^{r-1}(M)\right) \subseteq$ $C_{X}^{r-1}(M)$. Then the inclusion statement $C_{X}^{r-1}(M) \subseteq D(X+V)$ is [46, Proposition II.1.7], using [46, Definition II.1.6]. We let $\varphi \in C_{X}^{r-1}(M)$. It holds $X \varphi, \mathcal{L}_{\alpha, \phi_{\alpha}} \varphi \in C^{r-1}(M)$ since the flow is $C^{r}$. Recalling the weight $\phi_{\alpha}$ of the transfer operator in (2.18), with generating function $f \in C^{r-1}(M)$, we calculate and conclude:
$X \mathcal{L}_{\alpha, \phi_{\alpha}} \varphi=\left(X \phi_{\alpha}\right) \cdot \varphi \circ g_{-\alpha}+\phi_{\alpha} \cdot(X \varphi) \circ g_{-\alpha}=\left(f \circ g_{-\alpha}-f\right) \cdot \mathcal{L}_{\alpha, \phi_{\alpha}} \varphi+\phi_{\alpha} \cdot(X \varphi) \circ g_{-\alpha}$.

We set as the maximal spectral bound of the generator

$$
\begin{equation*}
\lambda_{\max }=\lambda_{\max }(s, t, q, p):=\left.\sup \Re \sigma(X+V)\right|_{W_{p}^{s, t, q}} . \tag{2.65}
\end{equation*}
$$

Lemma 2.4.10 (Discrete spectrum). Under the assumptions of Theorem 2.4.5 (including the choices for $p, s, t, q \in \mathbb{R}$ ), the set

$$
\left\{\left.\lambda \in \sigma(X+V)\right|_{W_{p}^{s, t, q}} \mid \Re \lambda>\lambda_{\text {min }}\right\}
$$

consists of isolated eigenvalues of finite multiplicity.
The discrete spectrum described in the previous lemma if $\lambda_{\max }>\lambda_{\min }$, is sometimes referred to as (Ruelle-Pollicott) resonances of $X+V$. In principle, the resonances depend on the choices $p, s, t$, and $q$ of the space $W_{p}^{s, t, q}$. We shall not enter into details here, but note that our main result in the next section shows that this dependence is mild, in particular, for the choice of $V$ there, $\lambda_{\max }$ is independent of $p, t, s$, and $q$.

Proof. Using Corollary 2.4.6, spectral radius of the resolvent is bounded from above by $\left|\Re z-\lambda_{\min }\right|^{-1}$. Assume $\left.\lambda \in \sigma(X+V)\right|_{W_{p}^{s, t, q}}$ such that $\Re \lambda>\lambda_{\text {min }}$. It follows from the Spectral Theorem for the Resolvent [46, Theorem V.1.13] that there exists $z \in \mathbb{C}$ (e.g. with $\Im z=\Im \lambda$ ) in the resolvent set of $X+V$ such that

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the spectral radius of the resolvent $\mathcal{R}_{z}$ has a lower bound given by

$$
|z-\lambda|^{-1}=(\Re z-\Re \lambda)^{-1}>\left(\Re z-\lambda_{\min }\right)^{-1} .
$$

Since $\Im \lambda$ was arbitrary we conclude.
The following notation associated to the eigenvalue spectrum is needed in Section 2.5 for the statement and proof of Theorem 2.5.7. We assume for the rest of this subsection

$$
\lambda_{\max }=\lambda_{\max }(s, t, q, p)>\lambda_{\min }(s, t, p)=\lambda_{\min },
$$

for any fixed choice $p \in[1, \infty]$ and $-s<0<q \leqslant t<r-1+s$. By Lemma 2.4.10 each $\left.\lambda \in \sigma(X+V)\right|_{W_{p}^{s, t, q}}$ such that $\Re \lambda>\lambda_{\text {min }}$ has a finite geometric multiplicity $n_{\lambda} \in \mathbb{N}$ and finite algebraic multiplicities $m_{\lambda, i} \in \mathbb{N}, 1 \leqslant i \leqslant n_{\lambda}$, with generalized eigenstates

$$
\mathcal{D}_{(\lambda, i, j)} \in D(X+V), \quad 1 \leqslant j \leqslant m_{\lambda, i},
$$

satisfying

$$
(X+V-\lambda)^{j} \mathcal{D}_{(\lambda, i, j)}=0 \quad \text { and } \quad \text { if } j>1:(X+V-\lambda)^{j-1} \mathcal{D}_{(\lambda, i, j)} \neq 0
$$

Moreover, to each geometric eigenvector there is associated a projector $\Pi_{\lambda, i}$ and a nil-potent operator $\mathcal{N}_{\lambda, i}$ of finite ranks such that

$$
\begin{align*}
& \Pi_{\lambda_{1}, i_{1}} \Pi_{\lambda_{2}, i_{2}} \equiv 0, \quad \mathcal{N}_{\lambda_{1}, i_{1}} \mathcal{N}_{\lambda_{2}, i_{2}} \equiv 0 \text { if } \quad \lambda_{1} \neq \lambda_{2} \quad \text { or } \quad i_{1} \neq i_{2},  \tag{2.66}\\
& \Pi_{\lambda_{1}, i_{1}} \mathcal{N}_{\lambda_{2}, i_{2}}=\mathcal{N}_{\lambda_{2}, i_{2}} \Pi_{\lambda_{1}, i_{1}}=\left\{\begin{array}{ll}
\mathcal{N}_{\lambda_{2}, i_{2}} & \text { if } \lambda_{1}=\lambda_{2} \text { and } i_{1}=i_{2} \\
0 & \text { if } \lambda_{1} \neq \lambda_{2} \text { or } i_{1} \neq i_{2}
\end{array},\right. \\
& \mathcal{N}_{\lambda, i}^{m_{\lambda, i}-1} \equiv 0 .
\end{align*}
$$

Note that the projector $\Pi_{\lambda, i}$ can be written as a finite rank operator

$$
\begin{equation*}
\Pi_{\lambda, i}=\sum_{j=1}^{n_{\lambda, i}} \mathcal{D}_{(\lambda, i, j)} \otimes \mathcal{O}_{(\lambda, i, j)}, \tag{2.67}
\end{equation*}
$$

where the dual vectors $\mathcal{O}_{(\lambda, i, j)} \in D(X+V)^{\prime}$ satisfy

$$
\mathcal{O}_{\left(\lambda_{1}, i_{1}, j_{1}\right)}\left(\mathcal{D}_{\left(\lambda_{2}, i_{2}, j_{2}\right)}\right)= \begin{cases}1, & \text { if }\left(\lambda_{1}, i_{1}, j_{1}\right)=\left(\lambda_{2}, i_{2}, j_{2}\right) \\ 0, & \text { otherwise } .\end{cases}
$$

We shall use the following Dolgopyat-type condition, adapted from [19, Assumption 3A], on the resolvent $\mathcal{R}_{z}=(z-X-V)^{-1}$, to control the remainder term $\mathcal{E}_{T, x}$ in (2.1) in Theorem 2.5.7 (to reduce to the case studied by Butterley, consider the renormalized semi-group $e^{-\lambda_{\max } \alpha} \mathcal{L}_{\alpha, \phi_{\alpha}}$ with generator $X+V-\lambda_{\max }$ and resolvent $\left.\mathcal{R}_{z+\lambda_{\max }}\right)^{4}$ :

Condition 2.4.11 (Spectral gap with (Dolgopyat) bounds). There exists

$$
\delta \in\left(\lambda_{\min }(s, t, p), \lambda_{\max }(s, t, q, p)\right)
$$

so that the following holds: For some $a>0, b>0, C>0$, some

$$
\gamma \in\left(0,1 / \log \left(1+\left(\lambda_{\max }-\delta\right) / a\right)\right),
$$

and for all $z \in \mathbb{C}$ with $\Re z=a$ and $|\Im z| \geqslant b$, we have

$$
\left\|\mathcal{R}_{z+\lambda_{\max }^{\tilde{n}}}\right\|_{W_{p}^{s, t, q}} \leqslant C^{\tilde{n}}\left|\Re z+\left(\lambda_{\max }-\delta\right)\right|^{-\tilde{n}}, \quad \text { where } \tilde{n}=\lceil\gamma \log |\Im z|\rceil
$$

It is well known that if $\left\|\mathcal{L}_{\alpha, \phi_{\alpha}}\right\|_{W_{p}^{s, t, q} \rightarrow W_{p}^{s, t, q}} \leqslant C e^{\lambda_{\max } \alpha}$ for all $\alpha$ and if $\mathcal{R}_{z}$ enjoys Lasota-Yorke estimates for $\lambda_{\min }(s, t)$ on $W_{p}^{s, t, q}$, in the sense of Theorem 2.4.5, then Condition 2.4.11 for some constant $\delta$ implies a spectral gap for the same $\delta$, in the sense that

$$
\begin{equation*}
\left.\sigma(X+V)\right|_{W_{p}^{s, t, q}} \cap\{\Re \lambda>\delta\} \text { is a finite set, } \tag{2.68}
\end{equation*}
$$

see e.g. [19, Theorem 1]. (Note that [19, Assumption 1] follows from the facts that $W_{p}^{s, t, q} \subset W_{p}^{s, t, q-1},\|(X+V) \varphi\|_{W_{p}^{s, t, q-1}} \leqslant C\|\varphi\|_{W_{p}^{s, t, q}}$ for some constant $C>0$, using Lemma 2.3.3 and Lemma 2.3.5, and

$$
\begin{equation*}
e^{-\lambda_{\max } \alpha} \mathcal{L}_{\widetilde{\alpha}, \phi_{\widetilde{\alpha}}} \varphi-\varphi=\left(X+V-\lambda_{\max }\right) \int_{0}^{\alpha} e^{-\lambda_{\max } \widetilde{\alpha}} \mathcal{L}_{\widetilde{\alpha}, \phi_{\tilde{\alpha}}} \varphi \mathrm{d} \widetilde{\alpha} \tag{2.69}
\end{equation*}
$$

for all $\varphi \in W_{p}^{s, t, q}$.)

[^4]Beware that even when $W_{p}^{s, t, q}$ is a Hilbert space, the operator $X+V$ is not selfadjoint a priori, so the existence of a spectral gap for $X+V$ with $\delta$ does not imply a spectral gap with bounds on the resolvent in general. (In the self-adjoint case, classical bounds on the iterated resolvent $\mathcal{R}_{z}^{n}$ in terms of the distance between $z$ and the spectrum give bounds stronger than Condition 2.4.11.)
See also Remark 2.5.11 for a further discussion of Condition 2.4.11.

### 2.4.4 Proof of Lemma 2.4.1

We need some preparations. We recall the quantities $\|F\|_{-, I},\|F\|_{+, I},\|F\|_{0, I}$ given below (2.47). We introduce an arrow relation as used by Baladi and Tsujii in [10, p.16].

Definition 2.4.12 (Arrow relation). Let $n, \ell \in \mathbb{Z} \geqslant 0$ and $\sigma, \tau \in\{-,+, 0\}$. We write

$$
\begin{aligned}
(\tau, \ell) \hookrightarrow_{I}(\sigma, n) & \Leftrightarrow\left\{\begin{array}{l}
\tau=+ \text { and } 2^{n-\ell} \leqslant 2^{4}\|F\|_{+, I} \\
\tau=\sigma=- \text { and } 2^{n-\ell} \geqslant 2^{-4}\|F\|_{-, I}
\end{array},\right. \\
\ell \hookrightarrow_{I}(\sigma, n) & \Leftrightarrow\left\{\begin{array}{l}
\tau=\sigma=0 \text { and } 2^{n-\ell} \leqslant 2^{4}\|F\|_{0, I}, \\
\sigma=- \text { and } \tau=0
\end{array}\right.
\end{aligned}
$$

and $(\tau, \ell) \leftrightarrow_{I}(\sigma, n)$ in the other cases.
We recall the function $c$ defined in (2.29). We let $c^{\prime}$ be analogously defined for $s^{\prime} \leqslant s, t^{\prime} \leqslant t, q^{\prime} \leqslant q$. We have for some constant $C>0$, for all fixed $\tau \in\{-,+, 0\}, \ell \in \mathbb{Z}_{\geqslant 0}$

$$
\begin{align*}
\sum_{(\tau, \ell) \hookrightarrow I(\sigma, n)} 2^{c(\sigma) n-c(\tau) \ell} & =\sum_{(\tau, \ell) \hookrightarrow_{I}(\sigma, n)} 2^{(c(\sigma)-c(\tau)) n+c(\tau)(n-\ell)} \leqslant \sum_{(\tau, \ell) \hookrightarrow I(\sigma, n)} 2^{c(\tau)(n-\ell)} \\
& \leqslant C \max \left\{\|F\|_{+, I}^{t},\|F\|_{-, I}^{s}\right\} . \tag{2.70}
\end{align*}
$$

An analogous estimate holds for all fixed $\sigma, n$. Similarly, we find either for all fixed $\ell$ or for all fixed $\sigma, n$

$$
\begin{equation*}
\sum_{\ell \leftrightarrows_{I}(\sigma, n)} 2^{c(\sigma) n-q \ell} \leqslant C \max \left\{1,\|F\|_{0, I}^{q}\right\} . \tag{2.71}
\end{equation*}
$$

We recall the norm of the Hilbert space $\ell_{2}^{c}$ (and analogously $\ell_{2}^{\prime}$ ) given in (2.30). Clearly, we have the inclusion $\ell_{2}^{c} \subseteq \ell_{2}^{c^{\prime}}$. We recall the definitions of $\Psi_{\sigma, n}^{\mathrm{Op}}$ in (2.23).
2.4 Properties of the transfer operator, the generator and its resolvent

We let given a family of pairwise disjoint sets

$$
\mathcal{I}:=\{I \subseteq K\} .
$$

For every $\left(a_{\tau, \ell}\right)=a \in L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{\prime}\right)$ we set

$$
\begin{align*}
& \left(Q_{\hookrightarrow}^{\mathrm{Op}} a\right)_{\sigma, n}:=\Psi_{\sigma, n}^{\mathrm{Op}} \sum_{I \in \mathcal{I}} \sum_{(\tau, \ell) \hookrightarrow_{I}(\sigma, n)} 1_{I I} a_{\tau, \ell}, \\
& \left(Q_{\hookrightarrow}^{\mathrm{Op}} a\right)_{\sigma, n}:=\Psi_{\sigma, n}^{\mathrm{Op}} \sum_{I \in \mathcal{I}} \sum_{l \hookrightarrow_{I}(\sigma, n)} 1_{\mid I} a_{0, l} . \tag{2.72}
\end{align*}
$$

Lemma 2.4.13 (Boundedness I). For all $p \in[1, \infty]$ the map

$$
Q_{\hookrightarrow}^{\mathrm{Op}}: L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{c}\right) \rightarrow L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{c}\right)
$$

is a bounded linear operator. Moreover, for some constant $C \geqslant 0$, for every $f \in L_{\infty}\left(\mathbb{R}^{d}, \mathbb{R} \backslash\{0\}\right)$ and every $a \in L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{c}\right)$, it holds

$$
\left\|Q_{\rightarrow I}^{\mathrm{Op}} a\right\|_{L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{c}\right)} \leqslant C \sup _{I \in \mathcal{I}} \max \left\{\|F\|_{+, I}^{t},\|F\|_{-, I}^{s}\right\}\left\|f_{\mid I}\right\|_{L_{\infty}}\left\|\sum_{I \in \mathcal{I}} \frac{1}{f_{\mid I}}\right\| a\left\|_{\ell_{2}}\right\|_{L_{p}} .
$$

Let $c^{\prime}(0)=c(0)$. Then for all $p \in[1, \infty]$ the map

$$
Q_{\hookrightarrow 0, I}^{\mathrm{Op}}: L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{c^{\prime}}\right) \rightarrow L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{c}\right)
$$

is a bounded linear operator. Moreover, for every $a \in L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{\prime}\right)$ it holds

$$
\left\|Q_{\rightarrow 0, \mathcal{I}}^{\mathrm{Op}} a\right\|_{L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{c}\right)} \leqslant C \sup _{I \in \mathcal{I}} \max \left\{1,\|F\|_{0, I}^{q}\right\}\left\|f_{\mid I}\right\|_{L_{\infty}}\left\|\sum_{I \in \mathcal{I}} \frac{1}{f_{\mid I}}\right\| a\left\|_{\ell_{2}^{c^{\prime}}}\right\|_{L_{p}} .
$$

Proof. For every $b \in \ell_{2}^{c}$ we set

$$
(Q b)_{\sigma, n}:=\Psi_{\sigma, n} b_{\sigma, n},
$$

and for every $a \in L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{c}\right)$ we set

$$
(\mathcal{K} a)_{\sigma, n}:=\sum_{I \in \mathcal{I}} \sum_{(\tau, \ell) \longrightarrow I(\sigma, n)} 1_{\mid I} a_{\tau, \ell} .
$$

2 Horocycle averages on closed manifolds

We let $Q^{\mathrm{Op}}$ be the operator in (2.27) associated to $Q$. We note that

$$
Q_{\hookrightarrow}^{\mathrm{Op}}=Q^{\mathrm{Op}} \circ \mathcal{K} .
$$

Using Lemma 2.3.1, we bound for some constant $C_{1}>0$

$$
\left\|Q^{\mathrm{Op}} \mathcal{K} a\right\|_{L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{c}\right)} \leqslant C_{1}\|\mathcal{K} a\|_{L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{c}\right)}
$$

We estimate with constants $C_{2}, C_{3}>0$, using pairwise disjointness of elements $I \in \mathcal{I}$, Cauchy-Schwarz and the bound in (2.70),

$$
\begin{aligned}
& \|\mathcal{K} a\|_{L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{c}\right)}=\left\|\left(\sum_{\sigma, n} 4^{c(\sigma) n}\left(\sum_{I \in \mathcal{I}(\tau, \ell) \rightarrow I(\sigma, n)} \sum_{\mid I} a_{\tau, \ell}\right)^{2}\right)^{\frac{1}{2}}\right\|_{L_{p}} \\
& \leqslant\left\|\left(\sum_{\sigma, n} \sum_{I \in \mathcal{I}}\left(\sum_{(\tau, \ell) \rightarrow{ }_{I}(\sigma, n)} 2^{c(\sigma) n-c(\tau) \ell}\right) \sum_{(\tau, \ell) \rightarrow{ }_{I}(\sigma, n)} 2^{c(\sigma) n+c(\tau) \ell}\left|1_{\mid I} a_{\tau, \ell}\right|^{2}\right)^{\frac{1}{2}}\right\|_{L_{p}} \\
& \leqslant C_{2}\left\|\left(\sum_{I \in \mathcal{I}} \max \left\{\|F\|_{+, I}^{t},\|F\|_{-, I}^{s}\right\} \sum_{\sigma, n} \sum_{(\tau, \ell) \rightarrow I(\sigma, n)} 2^{c(\sigma) n+c(\tau) \ell}\left|1_{\mid I} a_{\tau, \ell}\right|^{2}\right)^{\frac{1}{2}}\right\|_{L_{p}} \\
& =C_{2}\left\|\left(\sum_{I \in \mathcal{I}} \max \left\{\|F\|_{+, I}^{t},\|F\|_{-, I}^{s}\right\} \sum_{\tau, \ell} 2^{2 c(\tau) \ell}\left|1_{\mid I} a_{\tau, \ell}\right|^{2} \sum_{(\tau, \ell) \rightarrow I(\sigma, n)} 2^{c(\sigma) n-c(\tau) \ell}\right)^{\frac{1}{2}}\right\|_{L_{p}} \\
& \leqslant C_{3}\left\|\left(\sum_{I \in \mathcal{I}} \max \left\{\|F\|_{+, I}^{2 t},\|F\|_{-, I}^{2 s}\right\}\left\|f_{\mid I}\right\|_{L_{\infty}}^{2} \sum_{\tau, \ell} 2^{2 c(\tau) \ell}\left|\frac{1}{f_{\mid I}} a_{\tau, \ell}\right|^{2}\right)^{\frac{1}{2}}\right\|_{L_{\infty}} \\
& \leqslant C_{3} \sup _{I \in \mathcal{I}} \max \left\{\|F\|_{+, I}^{t},\|F\|_{-, I}^{s}\right\}\left\|f_{\mid I}\right\|_{L_{\infty}}\left\|\left(\sum_{I \in \mathcal{I}} \sum_{\tau, \ell} 2^{2 c(\tau) \ell}\left|\frac{1}{f_{\mid I}} a_{\tau, \ell}\right|^{2}\right)^{\frac{1}{2}}\right\|_{L_{p}} .
\end{aligned}
$$

The statement about $Q_{\rightarrow 0, 工}^{\mathrm{Op}}$ follows analogously, using (2.71).
We recall (see above (2.47)) that $F$ is assumed to be ( $\Theta^{\circ}, \Theta$ )-hyperbolic on $K$ and recall the maps $\tilde{\Phi}_{-}^{\circ}, \tilde{\Phi}_{0}^{\circ}$ assumed in Definition 2.2.2 which we use to construct
$\tilde{\Psi}_{\sigma, n}^{\circ}$ defined in (2.24). We set

$$
\mathcal{J}:=\begin{gather*}
\{(-, 0,0,0),(-, 0,+, 0),(0,0,+, 0),(+, 0,-, 0)\} \cup  \tag{2.73}\\
\{(+, \ell,+, 0),(+, \ell, 0,0),(0, \ell, 0,0) \mid \ell \geqslant 0\} \cup\{(-, 0,-, n) \mid n \geqslant 0\}
\end{gather*}
$$

We recall the arrow notation f from Definition 2.4.12.
Lemma 2.4.14 (Directional inequality). Let $(\tau, \ell) \leftrightarrow_{I}(n, \sigma)$ and $(\tau, \ell, \sigma, n) \notin$ J. Let $\eta \in \operatorname{supp} \tilde{\Psi}_{\tau, \ell}^{\circ}$ and $\xi \in \operatorname{supp} \Psi_{\sigma, n}$. Set

$$
m(\tau):= \begin{cases}\max \{n, \ell\}, & \text { if } \tau \in\{-, 0\}  \tag{2.74}\\ n, & \text { if } \tau=+\end{cases}
$$

Then, for some $C>0$ and for all $y \in I$ it holds

$$
\left|\left(\mathrm{D}_{y} F\right)^{\operatorname{tr}} \eta-\xi\right| \geqslant C 2^{m(\tau)} \min \left\{1,\|F\|_{-, I},\|F\|_{0, I}\right\}
$$

Proof. This can be seen case-by-case for admissible $\sigma, \tau$ as follows. We recall the set $\mathcal{J}$ defined in (2.73). We let $(\tau, \ell) \leftrightarrow_{I}(n, \sigma)$ such that $(\tau, \ell, \sigma, n) \notin \mathcal{J}$. Due to the construction of $\Psi_{\sigma, n}$ and $\tilde{\Psi}_{\tau, \ell}^{\circ}$, respectively, if $n \geqslant 1$ then $2^{n-1} \leqslant|\xi| \leqslant 2^{n+1}$ and if $\ell \geqslant 1$ then $2^{\ell-2} \leqslant|\eta| \leqslant 2^{\ell+2}$. We assume first $c(\sigma) \leqslant c(\tau)$. Let $\tau=+$. Then $2^{n-\ell}>2^{4}\|F\|_{+, I}$ and moreover, the exclusion of $\mathcal{J}$ implies $n \geqslant 1$. Using the triangle inequality, we find

$$
\left|\left(\mathrm{D}_{y} F\right)^{\operatorname{tr}} \eta-\xi\right| \geqslant|\xi|-\left|\left(\mathrm{D}_{y} F\right)^{\operatorname{tr}} \eta\right| \geqslant 2^{n-1}-\|F\|_{+, I} 2^{\ell+2} \geqslant 2^{n-1}-2^{n-2} \geqslant 2^{n-2}
$$

The case $\tau=0$ is analogous. Just note that we have also the estimate

$$
2^{n-1}-\|F\|_{0, I} 2^{\ell+2} \geqslant\|F\|_{0, I} 2^{l+3}-\|F\|_{0, I} 2^{\ell+2}
$$

If $\tau=-$ it holds $2^{n-\ell}<2^{-4}\|F\|_{-, I}$. The exclusion of $\mathcal{J}$ implies $l \geqslant 1$. Using the triangle inequality, we find
$\left|\left(\mathrm{D}_{y} F\right)^{\operatorname{tr}} \eta-\xi\right| \geqslant\left|\left(\mathrm{D}_{y} F\right)^{\operatorname{tr}} \eta\right|-|\xi| \geqslant\|F\|_{-, I} 2^{\ell-2}-2^{n+1}>2^{n+2}-2^{n+1} \geqslant 2^{n+1}$.

On the other hand we have also the estimate

$$
\|F\|_{-, I} 2^{\ell-2}-2^{n+1}>\|F\|_{-, I} 2^{\ell-2}-\|F\|_{-, I} 2^{l-3}
$$

Now we assume $c(\sigma)>c(\tau)$. We assume first $\tau=-$. Then $\sigma \in\{0,+\}$. We recall
that $F$ is cone-hyperbolic (see Definition 2.2.2). The exclusion of $\mathcal{J}$ implies $n \neq 0$ or $l \neq 0$. Together with the first compact inclusion in (2.13) we conclude that the angle between $\left(\mathrm{D}_{y} F\right)^{\operatorname{tr}} \eta$ and $\xi$ is bounded from below. This implies a lower bound $\geqslant C 2^{\max \{n, l\}}$ for the distance in both cases where $C>0$ is some constant. We assume now $\tau=0$ which implies $\sigma=+$. The reasoning is analogous as for $\tau=-$, using the second compact inclusion in (2.13) to bound the angle between $\left(\mathrm{D}_{y} F\right)^{\mathrm{tr}} \eta$ and $\xi$ from below.

Lemma 2.4.15. Let $p \in[1, \infty], b \in L_{p}$ and let $(\tau, \ell) \leftrightarrow_{I}(\sigma, n)$ and $(\tau, n, \sigma, \ell) \neq$ $\mathcal{J}$. It holds for the local transfer operator $L_{F, f_{\mid I}}$

$$
\left\|\Psi_{\sigma, n}^{\mathrm{Op}} L_{F, f_{I I}} \widetilde{\Psi}_{\tau, \ell}^{\circ \mathrm{Op}} b\right\|_{L_{p}} \leqslant C_{3}\left(F, f_{\mid I}\right) 2^{-(r-1) m(\tau)}\|b\|_{L_{p}}
$$

where for some $C \geqslant 1$ it holds
$C_{3}\left(F, f_{\mid I}\right) \leqslant C \max \left\{1,\|F\|_{-, I}^{1-r},\|F\|_{0, I}^{1-r}\right\}\|\mathrm{D} F\|_{C^{r-1}}^{k}\|f\|_{C^{r-1}} \sup _{y \in K}\left|\operatorname{det} \mathrm{D}_{y} F\right|^{-\frac{1}{p}}$.
Proof. This is analogous to the proof of Lemma 2.3.3, except that we have to deal with the additional composition operation by the map $F$. We set $f:=f_{\mid I}$. We expand the convolution and inverse Fourier transform

$$
\begin{aligned}
& \Psi_{\sigma, n}^{\mathrm{Op}} L_{F, f} \widetilde{\Psi}_{\tau, \ell}^{\circ} \mathrm{Op} \\
& \mathrm{O} \\
&=C \int_{\mathbb{R}^{4 d}} e^{\mathrm{i} \eta(F(z)-y)} e^{\mathrm{i} \xi(x-z)} f(z) b(y) \Psi_{\sigma, n}(\xi) \widetilde{\Psi}_{\tau, \ell}^{\circ}(\eta) \mathrm{d} \eta \mathrm{~d} \xi \mathrm{~d} z \mathrm{~d} y \\
&=C \int_{\mathbb{R}^{d}} V_{\sigma, n}^{\tau, \ell}(x, y) b(F(y))|\operatorname{det} \mathrm{D} F(y)| \mathrm{d} y,
\end{aligned}
$$

for some constant $C>0$ and where we set

$$
\begin{equation*}
V_{\sigma, n}^{\tau, \ell}(x, y):=\int_{\mathbb{R}^{3 d}} e^{-\mathrm{i} \eta F(y)} e^{\mathrm{i} \xi x} \Psi_{\sigma, n}(\xi) \widetilde{\Psi}_{\tau, \ell}^{\circ}(\eta) e^{\mathrm{i}(\eta F(z)-\xi z)} f(z) \mathrm{d} z \mathrm{~d} \eta \mathrm{~d} \xi . \tag{2.75}
\end{equation*}
$$

We transform (2.75), first integrating by parts $\lfloor r\rfloor-1$-times in $z$ (see Lemma B. 3 with function $G(z):=\eta F(z)-\xi z$ which has a gradient bounded from below by Lemma 2.4.14). Therefore we replace $f(z)$ in (2.75) with another function $V_{[r\rfloor-1}(z, \eta, \xi)$ which satisfies the iterative construction given in Lemma B. 3 (B.1). Using Lemma 2.4.14 and Lemma B. 3 (B.2), we estimate for some constant $C \geqslant 1$

$$
\left\|V_{\lfloor r-1\rfloor}\right\|_{C^{0}} \leqslant C \widetilde{C}_{1} 2^{-m(\tau)\lfloor r-1\rfloor} \max \left\{1,\|F\|_{-, I}^{-\lfloor r-1\rfloor},\|F\|_{0, I}^{-\lfloor r-1\rfloor}\right\}\|f\|_{C^{[r-1\rfloor}},
$$

where $\widetilde{C}_{1}:=\sup _{(z, \eta, \xi) \in \text { supp } f} \max _{0 \leqslant|\gamma| \leqslant\lfloor r-1\rfloor}| |\left(\mathrm{D}_{z} F\right)^{\operatorname{tr}} \eta-\xi \left\lvert\, \partial_{z}^{\gamma} \frac{\left(\mathrm{D}_{z} F F\right)^{\mathrm{tr}} \eta-\xi}{\left|\left(\mathrm{D}_{z} F\right)^{\operatorname{tr}} \eta-\xi\right|^{2} \mid}{ }^{\lfloor r-1\rfloor}\right.$. Moreover, this function is a $C^{\tilde{r}}$-map for $\tilde{r}:=r-\lfloor r\rfloor$. Using Lemma B. 5 (in there we take $\epsilon=L^{-1}=2^{-m(\tau)}$ ), we proceed with a regularized integration by parts
in $z$. This yields

$$
\begin{equation*}
V_{\sigma, n}^{\tau, \ell}(x, y)=\int_{\mathbb{R}^{3 d}} e^{\mathrm{i} \eta(F(z)-F(y))} e^{\mathrm{i} \xi(x-z)} \Psi_{\sigma, n}(\xi) \tilde{\Psi}_{\tau, \ell}^{\circ}(\eta) V_{r-1}(z, \eta, \xi) \mathrm{d} z \mathrm{~d} \eta \mathrm{~d} \xi, \tag{2.76}
\end{equation*}
$$

where $V_{r-1}$ is given in (B.4) in Lemma B. 5 with bound

$$
\begin{equation*}
\left\|V_{r-1}\right\|_{C^{0}} \leqslant C \widetilde{C}_{2} 2^{-m(\tau) \tilde{r}} \max \left\{1,\|F\|_{-, I}^{\|_{-,}},\|F\|_{0, I}^{-\widetilde{r}}\right\} \tag{2.77}
\end{equation*}
$$

where $\tilde{C}_{2}:=\sup _{(z, \eta, \xi) \in \operatorname{supp} f}\left(1+\left|\left(\left(\mathrm{D}_{z} F\right)^{\operatorname{tr}} \eta-\xi\right)\right|\right)\left\|\frac{\left(\left(\mathrm{D}_{z} F\right)^{\mathrm{tr}} \eta-\xi\right) V_{[r-1]}(\cdot, \eta, \xi)}{\left|\left(\mathrm{D}_{z} F\right)^{\mathrm{tr} r} \eta-\xi\right|^{2}}\right\|_{C^{\dot{r}}}$.
We now substitute $\xi \rightarrow 2^{\ell} \xi^{\prime}$ and $\eta \rightarrow 2^{n} \eta^{\prime}$ in (2.76). By construction the function $V_{r-1}\left(z, 2^{n} \eta^{\prime}, 2^{\ell} \xi^{\prime}\right)$ is uniformly bounded in $n$ and $\ell$ in the $C^{\infty}$-norm with respect to $\eta^{\prime}$ and $\xi^{\prime}$. We transform (2.76), integrating by parts $d+1$-times in $\xi^{\prime}$ if $\left|2^{n}(z-x)\right|>1$, and $d+1$-times in $\eta^{\prime}$ if $\left|2^{\ell}(F(z)-F(y))\right|>1$, which yields for some constant $C_{1}>0$

$$
V_{\sigma, n}^{\tau, \ell}(x, y)=C_{1} \int_{\mathbb{R}^{3} d} \frac{e^{\mathrm{i} 2^{\ell} \eta^{\prime}(F(z)-F(y))}}{u_{\ell}(F(z)-F(y))} \frac{e^{\mathrm{i} 2^{n^{\prime}} \xi^{\prime}(x-z)}}{u_{n}(x-z)} 2^{d(\ell+n)} \tilde{V}_{\sigma, n}^{\tau, \ell}\left(z, \eta^{\prime}, \xi^{\prime}\right) \mathrm{d} z \mathrm{~d} \eta^{\prime} \mathrm{d} \xi^{\prime},
$$

where $\tilde{V}_{\sigma, n}^{\tau, \ell}\left(z, \eta^{\prime}, \xi^{\prime}\right)$ together with $u_{n}: \mathbb{R}^{d} \rightarrow(0,1]: x \mapsto \begin{cases}1 & , \text { if }\left|2^{n} x\right| \leqslant 1 \\ \left|2^{n} x\right|^{d+1} & , \text { else }\end{cases}$ replaces $\Psi_{\sigma, n}\left(2^{\ell} \xi^{\prime}\right) \widetilde{\Psi}_{\tau, \ell}^{\circ}\left(2^{n} \eta^{\prime}\right) V_{r-1}\left(z, 2^{n} \eta^{\prime}, 2^{\ell} \xi^{\prime}\right)$ in (2.76). Since we only derived $V_{r-1}$ with respect to $\eta^{\prime}$ and $\xi^{\prime}$, respectively, the $C^{0}$-norm of $\tilde{V}_{\sigma, n}^{\tau, \ell}\left(z, \eta^{\prime}, \xi^{\prime}\right)$ is controlled by the upper bound given in (2.77). We recall that $\xi^{\prime}, \eta^{\prime}$ are uniformly bounded. We estimate trivially for some constant $C_{2} \geqslant 1$

$$
\begin{aligned}
& \| \Psi_{\sigma, n}^{\mathrm{Op}} L_{F, f} \tilde{\Psi}_{\tau, \ell}^{\circ} \mathrm{Op} \\
& \|_{L_{p}} \leqslant C\left\|\tilde{V}_{\sigma, n}^{\tau, \ell}\right\|_{C^{0}} 2^{d(n+\ell)}\left\|\frac{1}{u_{n}} *\left(\left(b * \frac{1}{u_{\ell}}\right) \circ F\right)\right\|_{L_{p}} \\
& \leqslant C_{2} C_{3}(F, f) 2^{-m(\tau)(r-1)}\|b\|_{L_{p}},
\end{aligned}
$$

where we used twice Young's inequality in the last step.

We set for all $n \in \mathbb{Z}_{\geqslant 0}$, for all $\sigma \in\{-,+, 0\}$ and for all $I \in \mathcal{I}$

$$
\begin{equation*}
\left(Q_{\leftrightarrow, I}^{\mathrm{Op}} a\right)_{\sigma, n}:=\Psi_{\sigma, n}^{\mathrm{Op}} \sum_{(\tau, \ell) \nrightarrow I}(\sigma, n), \quad a_{\tau, \ell}, \quad\left(Q^{\mathrm{Op}} a\right)_{\sigma, n}:=\widetilde{\Psi}_{\sigma, n}^{\circ} \mathrm{O}_{\mathrm{p}} a_{\sigma, n} \tag{2.78}
\end{equation*}
$$

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Lemma 2.4.16 (Boundedness II). Let $c(+)-c^{\prime}(-)<r-1$. Then for all $p \in[1, \infty]$ the map $Q_{\leftrightarrow, I}^{\mathrm{Op}} L_{F, f_{I}} Q^{\mathrm{Op}}: L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{c^{\prime}}\right) \rightarrow L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{c}\right)$ is a bounded linear operator. In particular, it holds

$$
\left\|Q_{\leftrightarrow, I}^{\mathrm{Op}} L_{F, f_{I}} Q^{\mathrm{Op}}\right\|_{L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{\prime}\right) \rightarrow L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{c}\right)} \leqslant C C_{4}\left(F, f_{\mid I}\right),
$$

where for some $C \geqslant 1$ and some $k \geqslant 0$
$C_{4}(F, f) \leqslant C \max \left\{1,\|F\|_{-, I}^{1-r},\|F\|_{0, I}^{1-r}\right\} \max \left\{1,\|\mathrm{D} F\|_{C^{r-1}}^{k}\right\}\|f\|_{C^{r-1}} \sup _{y \in K}\left|\operatorname{det} \mathrm{D}_{y} F\right|^{-\frac{1}{p}}$
Proof. Let $\left(a_{\tau, \ell}\right)=a \in L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{c^{\prime}}\right)$. We have

$$
\begin{align*}
\left.\left\|Q_{\leftrightarrow \rightarrow, I}^{\mathrm{Op}} L_{F, f_{I}} Q^{\mathrm{Op}} a\right\|_{L_{p}\left(\mathbb{R}^{d}, \ell \ell\right.}\right) & =\left\|\left(\left.\sum_{\sigma, n} 4^{c(\sigma) n} \sum_{(\tau, \ell) \leftrightarrow \mapsto_{I}(\sigma, n)} \Psi_{\sigma, n}^{\mathrm{Op}} L_{F, f_{\mid I}} \widetilde{\Psi}_{\tau, \ell}^{\circ \mathrm{Op}} a_{\tau, \ell}\right|^{2}\right)^{\frac{1}{2}}\right\|_{L_{p}} \\
& \leqslant \sum_{(\tau, \ell) \nmid I(\sigma, n)} 2^{c(\sigma) n}\left\|\Psi_{\sigma, n}^{\mathrm{Op}} L_{F, f_{I}} \widetilde{\Psi}_{\tau, \ell}^{\circ \mathrm{Op}} a_{\tau, \ell}\right\|_{L_{p}} . \tag{2.79}
\end{align*}
$$

We recall the set of indices $\mathcal{J}$ in (2.73). We assume $(\tau, l, \sigma, n) \in \mathcal{J}$. Now we make three distinctions in the estimate of the corresponding part of the sum in (2.79). If $\tau \in\{+, 0\}$ then $n=0$ and $l \geqslant 0$. Then, using Young's inequality, for some $C \geqslant 1$

$$
\begin{aligned}
\sum_{l=0}^{\infty}\left\|\Psi_{\sigma, 0}^{\mathrm{Op}} L_{F, f_{I I}} \tilde{\Psi}_{\tau, \ell}^{\circ \mathrm{Op}} a_{\tau, \ell}\right\|_{L_{p}} & \leqslant C \sup _{z \in I}|f(z)|\left|\operatorname{det} \mathrm{D}_{z} F\right|^{-\frac{1}{p}} \sum_{l=0}^{\infty} 2^{-c^{\prime}(\tau) l} \sup _{l} 2^{2^{\prime}(\tau) l}\left\|a_{\tau, \ell}\right\|_{L_{p}} \\
& \leqslant C \sup _{z \in I}|f(z)|\left|\operatorname{det} \mathrm{D}_{z} F\right|^{-\frac{1}{p}}\|a\|_{L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{\prime}\right)}
\end{aligned}
$$

If $\tau=\sigma=-$ then $n \geqslant 0$ and $l=0$. Recall that $s<0$. Then, using Young's inequality,

$$
\sum_{n=0}^{\infty} 2^{s n}\left\|\Psi_{\sigma, n}^{\mathrm{Op}} L_{F, f_{\mid I}} \tilde{\Psi}_{-, 0}^{\prime \mathrm{Op}} a_{-, 0}\right\|_{L_{p}} \leqslant C \sup _{z \in I}|f(z)|\left|\operatorname{det} \mathrm{D}_{z} F\right|^{-\frac{1}{p}}\|a\|_{L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{\prime}\right)} .
$$

In the three remaining cases $n=l=0$ we estimate analogously, using Young's inequality. Now we assume $(\tau, l, \sigma, n) \notin \mathcal{J}$. We recall $m(\tau)$ defined in (2.74) in Lemma 2.4.14 and the constant $C_{3}(F, f)$ in Lemma 2.4.15. Using Lemma 2.4.15,
we estimate the remaining part of the sum in (2.79)

$$
\begin{align*}
& \sum_{\substack{(\tau, \ell) \leftrightarrow I(\sigma, n) \\
(\tau, l, \sigma, n) \notin \mathcal{J}}} 2^{c(\sigma) n} \| \Psi_{\sigma, n}^{\mathrm{Op}} L_{F, f_{I I}} \widetilde{\Psi}_{\tau, \ell}^{\circ} \mathrm{O} \mathrm{p} \\
& a_{\tau, \ell} \|_{L_{p}}  \tag{2.80}\\
& \leqslant C C_{3}(F, f) \sum_{\substack{(\tau, \ell) \nmid I, I(\sigma, n) \\
(\tau, l, \sigma, n) \notin \mathcal{J}}} 2^{c(\sigma) n-c^{\prime}(\tau) l-m(\tau)(r-1)} \sup _{\tau, \ell} 2^{c^{\prime}(\tau) l}\left\|a_{\tau, \ell}\right\|_{L_{p}}
\end{align*}
$$

$$
\leqslant C C_{3}(F, f)\|a\|_{L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{c_{2}^{\prime}}\right)}
$$

where the sums in $n, l$, respectively, in the right-hand side in (2.80) are bounded by geometric sums, using the assumption $c^{\prime}(-)>r-1-c(+)$. In particular, we find for (the worst-case since $0<c(0) \leqslant c(+)<r-1) \tau=-, \sigma=+$, if $l \geqslant n$ for all small enough $\epsilon>0$

$$
\begin{aligned}
(c(+)+\epsilon) n-c^{\prime}(-) l-m(-)(r-1) & \leqslant\left(c(+)+\epsilon-c^{\prime}(-)\right) l-m(-)(r-1) \\
& =\left(c(+)+\epsilon-c^{\prime}(-)-r+1\right) l<0,
\end{aligned}
$$

and an analogous estimate holds for $l<n$. We note

$$
\sup _{z \in I}|f(z)|\left|\operatorname{det} \mathrm{D}_{z} F\right|^{-\frac{1}{p}} \leqslant C\|f\|_{C^{r-1}} \sup _{z \in I}\left|\operatorname{det} \mathrm{D}_{z} F\right|^{-\frac{1}{p}} .
$$

We set $C_{4}\left(F, f_{\mid I}\right):=C C_{3}\left(F, f_{\mid I}\right)$. Combining the estimates for all the parts of the sum (2.79), we conclude.

Proof of Lemma 2.4.1. Let $s, q, t, p$ satisfy the hypotheses in Lemma 2.4.1. That is $s^{\prime}<s<0<q \leqslant t<r-1+s^{\prime}, q^{\prime}<q, t^{\prime}<t$ and $p \in[1, \infty]$. We put $c(-):=s$, $c(+):=t, c(0):=c^{\prime}(0):=q$ and $c^{\prime}(-):=c^{\prime \prime}(-):=s^{\prime}, c^{\prime}(+):=c^{\prime \prime}(+):=t^{\prime}$ and $c^{\prime \prime}(0):=q^{\prime}$. Then $c, c^{\prime}, c^{\prime \prime}$ satisfy (2.29), respectively, while $c, c^{\prime}$ satisfy the hypotheses in Lemma 2.4.13, and $c, c^{\prime \prime}$ that of Lemma 2.4.16. Let $\varphi \in W_{p, \Theta, F(K)}^{s, t, q}$. We set

$$
a_{\tau, \ell}:=L_{F, f} \Psi_{\tau, \ell}^{\circ \mathrm{Op}_{\mathrm{p}}} \varphi
$$

We have $a \in L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{c}\right) \subseteq L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{\prime}\right)$ because

$$
\begin{align*}
& \|a\|_{L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{\ell}\right)}=\left\|\left(\sum_{\tau, \ell} 4^{c(\tau) \ell}\left|a_{\tau, \ell}\right|^{2}\right)^{\frac{1}{2}}\right\|_{L_{p}}=\left\|\left(\sum_{\tau, \ell} 4^{c(\tau) \ell}\left|f \cdot\left(\Psi_{\tau, \ell}^{\circ \mathrm{O}_{\mathrm{p}}} \varphi\right) \circ F\right|^{2}\right)^{\frac{1}{2}}\right\|_{L_{p}} \\
& \leqslant\left\|f|\operatorname{det} \mathrm{D} F|^{-\frac{1}{p}}\right\|_{L_{\infty}(K)}\left\|\left(\sum_{\tau, \ell} 4^{c(\tau) \ell}\left|\Psi_{\tau, \ell}^{\circ \mathrm{O} \mathrm{p}} \varphi\right|^{2}\right)^{\frac{1}{2}}\right\|_{L_{p}(F(K))}  \tag{2.81}\\
& \leqslant\left\|f|\operatorname{det} \mathrm{D} F|^{-\frac{1}{p}}\right\|_{L_{\infty}}\|\varphi\|_{W_{p, \Theta^{s}, F(K)}^{s, t, q}} .
\end{align*}
$$

We set $b_{\tau, \ell}:=\Psi_{\tau, \ell}^{\circ \mathrm{Op}} \varphi$. By the first statement in Lemma 2.3.4, it holds $\varphi \in$ $W_{p, \Theta^{\circ}, F(K)}^{s^{\prime}, t^{\prime}, q^{\prime}}$ hence $\left(b_{\tau, \ell} \mid \tau \in\{-,+, 0\}, \ell \in \mathbb{Z}_{\geqslant 0}\right)=: b \in L_{p}\left(\mathbb{R}^{d}, \ell_{2}^{\prime^{\prime \prime}}\right)$. By assumption on $K$, we can decompose $K=\sqcup \mathcal{W}$ into finitely many open sets $\mathcal{W}$. For each component $\mathcal{W}$ we set

$$
a_{\mid \mathcal{W}, \tau, \ell}:=L_{F, f \mid \mathcal{W}} \Psi_{\tau, \ell}^{\circ \mathrm{O}_{\mathrm{p}}} \varphi \quad \text { and } \quad a:=\sum_{\mathcal{W}} a_{\mid \mathcal{W}} .
$$

By construction (see above (2.78)), it holds $\tilde{\Psi}_{\tau, \ell \mid \operatorname{supp} \Psi_{\tau, \ell}^{\prime}}^{\circ} \equiv 1$ hence $\Psi_{\tau, \ell}^{\circ \mathrm{Op}_{\mathrm{p}}} \varphi=$ $\tilde{\Psi}_{\tau, \ell}^{\circ}{ }^{\mathrm{Op}} \Psi_{\tau, \ell}^{\circ}{ }^{\mathrm{Op}} \varphi$. For each $\mathcal{W}$ there is a corresponding arrow relation given by Definition 2.4.12 and the restriction $f_{\mid \mathcal{W}}$ is also $C^{r-1}$. We rewrite

$$
\begin{align*}
\Psi_{\sigma, n}^{\mathrm{Op}} L_{F, f \mid \mathcal{W}} \varphi=\Psi_{\sigma, n}^{\mathrm{Op}} \sum_{\tau, \ell} a_{\mid \mathcal{W}, \tau, \ell}= & \Psi_{\sigma, n}^{\mathrm{Op}} \sum_{(\tau, \ell) \leftrightarrow \mathcal{W}(\sigma, n)} a_{\mid \mathcal{W}, \tau, \ell}+\Psi_{\sigma, n}^{\mathrm{Op}} \sum_{l \hookrightarrow \mathcal{W}(\sigma, n)} a_{\mid \mathcal{W}, 0, \ell} \\
& +\Psi_{\sigma, n}^{\mathrm{Op}} \sum_{(\tau, \ell) \leftrightarrow \mathcal{W}^{( }(\sigma, n)} L_{F, f_{\mid \mathcal{W}}} \widetilde{\Psi}_{\tau, \ell}^{\circ \mathrm{Op}} \Psi_{\tau, \ell}^{\circ \mathrm{Op}} \varphi . \tag{2.82}
\end{align*}
$$

We recall the definitions of the operators $Q_{\hookrightarrow,\{\mathcal{W}\}}^{\mathrm{Op}}, Q_{\rightarrow 0,\{\mathcal{W}\}}^{\mathrm{Op}}, Q_{\leftrightarrow}^{\mathrm{Op}}, \mathcal{W}$ given in (2.72) and in (2.78), respectively (in which we take $\mathcal{I}=\{\mathcal{W}\}$ and $I=\mathcal{W}$ ). We estimate, using the decomposition given in (2.82),

$$
\begin{aligned}
& \left\|L_{F, f} \varphi\right\|_{W_{p, \theta, K}^{s, t, q}}=\left\|\sum_{\mathcal{W}} L_{F, f \mid \mathcal{W}} \varphi\right\|_{W_{p, \theta, K}^{s, t, q}}=\left\|\left(\sum_{\sigma, n} 4^{c(\sigma) n}\left|\Psi_{\sigma, n}^{\mathrm{Op}} \sum_{\mathcal{W}} L_{F, f \mid \mathcal{W}} \varphi\right|^{2}\right)^{\frac{1}{2}}\right\|_{L_{p}} \\
& \leqslant 5\left\|Q_{\hookrightarrow,\{\mathcal{W}\}}^{\mathrm{Op}} a\right\|_{L_{p}\left(\mathbb{R}, \ell_{2}^{\mathrm{c}}\right)}+5\left\|Q_{\hookrightarrow 0,\{\mathcal{W}\}}^{\mathrm{Op}} a\right\|_{L_{p}\left(\mathbb{R}, \ell_{2}^{\mathrm{c}}\right)}+5 \sum_{\mathcal{W}}\left\|Q_{\rightarrow, \mathcal{W}}^{\mathrm{Op}} L_{F, f \mid \mathcal{W}} Q^{\mathrm{Op}} b\right\|_{L_{p}\left(\mathbb{R}, \ell_{2}^{c}\right)} .
\end{aligned}
$$

We conclude, using Lemma 2.4.13 and Lemma 2.4.16 together with the estimate given in (2.81).

### 2.5 Asymptotics of horocycle averages

In this section, we assume $r \geqslant 2$ and topological mixing of the Anosov flow $g_{\alpha}$. (Contact Anosov flows are topologically mixing [41, Theorem 3.6] and hence serve as examples for such Anosov flows $g_{\alpha}$ in the case $d=3$.) In order to define the horocycle flow in Definition 2.5.1 below we assume that the stable dimension $d_{-}=1$ and that the strong-stable distribution $E_{-}$is orientable. The stable manifolds of $M$ with respect to the flow $g_{\alpha}$ are those (non-compact) Riemannian submanifolds which are tangent to $E_{-}$. As consequence of topological mixing, each of those stable manifolds is dense in $M$ [49, p. 84].

### 2.5.1 Horocycle flows and integrals and main results (Theorem 2.5.7)

Definition 2.5.1 (Horocycle flow). A flow $h_{\rho}: M \rightarrow M$ in $\rho \in \mathbb{R}$ is called a stable horocycle flow if and only if for all $\rho \in \mathbb{R}$

$$
\partial_{\rho} h_{\rho} \in E_{-} \backslash\{0\} .
$$

Remark 2.5.2 (Unit speed parametrization). By the Stable Manifold Theorem (see e.g. [42, Theorem 8.12]), there exists a parametrization of stable manifolds by the arc-length induced by the Riemannian metric on $M$. Since we assumed that $E_{-}$is orientable, this yields the unit speed parametrization of the horocycle flow (i.e. $\left|\partial_{\rho} h_{\rho}\right| \equiv 1$ ).

Our main result, Theorem 2.5.7 provides a decomposition giving the $T$-asymptotics of the following horocycle integral:

Definition 2.5.3 (Horocycle integral). For all $\varphi \in C_{X}^{r-1}(M)$, for all $x \in M$ let

$$
\begin{equation*}
\gamma_{x}(\varphi, T):=\int_{0}^{T} \varphi \circ h_{\rho}(x) \mathrm{d} \rho . \tag{2.83}
\end{equation*}
$$

denote the horocycle integral of the horocycle flow $h_{\rho}$ for the observable $\varphi$ at base point $x$.

In Theorem 2.5.7 we reveal its connection to the eigendistributions of a weighted transfer operator for the Anosov flow $g_{-\alpha}$ introduced in Section 2.3.1, using renormalization dynamics to connect the stable flow with the Anosov flow. Results can be obtained for an unstable horocycle flow in an analogous way.

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Definition 2.5.4 (Pointwise renormalization time). A map $\tau: \mathbb{R}^{2} \times M \rightarrow \mathbb{R}$ which satisfies

$$
\begin{equation*}
g_{\alpha} \circ h_{\rho}(x)=h_{\tau(\rho, \alpha, x)} \circ g_{\alpha}(x), \quad \forall \rho, \alpha \in \mathbb{R}, \forall x \in M, \tag{2.84}
\end{equation*}
$$

is called a pointwise renormalization time for the stable flow $h_{\rho}$.
Remark 2.5.5. This definition of the renormalization time $\tau$ is the same as used by Marcus (denoted by s* in his notation) in [49, p.83] to study ergodic properties of the horocycle flow.

Lemma 2.5.6 (Existence and uniqueness). A pointwise renormalization time exists and is unique.

Proof. For every $x \in M$ and for every $\rho, \alpha \in \mathbb{R}$ we set $h_{\alpha, \rho}(x):=g_{\alpha} \circ h_{\rho} \circ g_{-\alpha}(x)$. By Definition 2.5.1 and the invariant splitting (2.2), we find $\partial_{\rho} h_{\alpha, \rho} \in E_{-, x} \backslash\{0\}$. Hence $h_{\alpha, \rho}(x)$ parametrizes the same stable manifold with respect to $\rho$ as $h_{\rho}(x)$. If there were two different pointwise renormalization times $\tau$, there would be $\rho_{1}<\rho_{2} \in \mathbb{R}$ such that $h_{\alpha, \rho}(x)=h_{\rho_{1}}(x)=h_{\rho_{2}}(x)$. By density of stable leaves and non-singularity of the flow $h_{\rho}$, there are no periodic points of $h_{\rho}$ hence $\rho_{1}=\rho_{2}$.

Further properties of the renormalization time $\tau$ are given in Proposition 2.5.13 below. Assuming $\partial_{\rho} \tau(0,-\alpha, \cdot) \in C^{r-1}(M)$ for all $\alpha \geqslant 0$, we will consider the potential $V$ defined by

$$
\begin{equation*}
V \equiv-\partial_{\alpha} \partial_{\rho} \tau(0,0, \cdot) \tag{2.85}
\end{equation*}
$$

Then $\phi_{\alpha}$ defined in (2.18) is just

$$
\begin{equation*}
\phi_{\alpha}:=\partial_{\rho} \tau(0,-\alpha, \cdot) . \tag{2.86}
\end{equation*}
$$

It follows from (iv) in Lemma 2.5.18 below that for any $p \in[1, \infty], t-s<r-1$ and $s<0<q \leqslant t$ the spectral bound $\lambda_{\max }=\left.\sup \Re \sigma(X+V)\right|_{W_{p}^{s, t, q}}$ for the generator satisfies

$$
\begin{equation*}
\lambda_{\max }=h_{\mathrm{top}} \tag{2.87}
\end{equation*}
$$

In the special case of unit speed horocycle flow (see Remark 2.5.2) it holds (using Proposition 2.5.13 (viii) below)

$$
\phi_{\alpha}=\operatorname{det} \mathrm{D} g_{-\alpha \mid E^{-}} .
$$

Hence if the strong stable distribution $E_{-}$is $C^{1}$ (see Proposition 2.5.10 where this holds true if $d=3$ under the contact assumption) and $r \geqslant 2$ then we find $\partial_{\rho} \tau(0,-\alpha, \cdot) \in C_{X}^{1}(M)$. In particular, our results apply to all $C^{1}$ time reparametrizations of the unit speed horocycle flow $h_{\rho}$ (this is analogous to [30, Remark 2.4]).
The following theorem will be proved at the end of Section 2.5.3:
Theorem 2.5.7 (Expansion of horocycle integrals). Let $g_{\alpha}$ be a topologically mixing $C^{r}$-Anosov flow, with $r \geqslant 2$, such that $E_{-}$is orientable and $d_{-}=1$. Let $\mu$ be the unique Borel measure which is invariant by the horocycle flow $h_{\rho}$. Assume for all $\alpha \geqslant 0$

$$
\phi_{\alpha}:=\partial_{\rho} \tau(0,-\alpha, \cdot) \in C^{r-1}(M) .
$$

Assume further that there exist $p \in[1, \infty]$, and $s<0<q \leqslant t$ with $t-s<r-1$ such that, for the corresponding anisotropic space $W_{p}^{s, t, q}$ it holds $\lambda_{\min }<\lambda_{\max }=$ $h_{\text {top }}$, with $\lambda_{\text {min }}=\lambda_{\min }(t, s, p)$ from (2.56). Then, for all $x \in M$ and $T \geqslant 0$ there exist, for each $\left.\lambda \in \sigma(X+V)\right|_{W_{p}^{s, t, q}}$ with $\Re \lambda>\lambda_{\text {min }}$, constants $c_{(\lambda, i, j)}(T, x) \in \mathbb{C}$ with

$$
\sup _{T>0, x \in M}\left|c_{(\lambda, i, j)}(T, x)\right|<\infty, \forall 1 \leqslant i \leqslant n_{\lambda}, 1 \leqslant j \leqslant m_{\lambda, i},
$$

such that, for any $\delta \in \mathbb{R}$ with

$$
\max \left\{\lambda_{\text {min }}, 0\right\} \leqslant \delta<h_{\text {top }}
$$

and any finite ${ }^{5}$ subset $\Lambda_{\delta}$ of

$$
\Sigma_{\delta}:=\left.\sigma(X+V)\right|_{W_{p}^{s, t, q}} \cap\{\lambda \in \mathbb{C} \mid \Re \lambda>\delta\},
$$

such that for all $\varphi \in C_{X}^{r-1}(M)$ and all $T \geqslant e$

$$
\begin{aligned}
& \int_{0}^{T} \varphi \circ h_{\rho}(x) \mathrm{d} \rho=\gamma_{x}\left(\mathcal{D}_{\left(h_{\text {top }, 1,1)}, T\right)} \mu(\varphi)\right. \\
& \quad+\sum_{\substack{\lambda \in \Lambda_{\delta} \\
\Re \lambda<h_{\text {top }}}} \sum_{1 \leqslant j \leqslant j \leqslant m_{\lambda, i}} T^{\frac{\lambda}{h_{\text {top }}}}(\log T)^{j-1} c_{(\lambda, i, j)}(T, x) \mathcal{O}_{(\lambda, i, j)}(\varphi)+\mathcal{E}_{T, x, \Lambda_{\delta}}(\varphi),
\end{aligned}
$$

where the dual eigendistributions $\mathcal{O}_{(\lambda, i, j)} \in D(X+V)^{\prime}$ are associated to the

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eigenvalue $\lambda$ by Lemma 2.4.10 (see (2.67)), and where

$$
\lim _{T \rightarrow \infty} \frac{\gamma_{x}\left(\mathcal{D}_{\left(h_{\text {top },}, 1,1\right)}, T\right)}{T}=1 \quad \text { and } \quad \lim _{T \rightarrow \infty} \frac{\mathcal{E}_{T, x, \Lambda_{\delta}}(\varphi)}{T}=0 .
$$

Moreover, if for some $c>0$ and some constant $C_{1}=C_{1}\left(\varphi, c, \Lambda_{\delta}\right) \geqslant 0$ for all $\alpha \geqslant 0$

$$
\begin{equation*}
\left\|\mathcal{L}_{\alpha, \phi_{\alpha}}\left(\sum_{\lambda \in \Lambda_{\delta}} \sum_{i=1}^{n_{\lambda}} \Pi_{\lambda, i} \varphi-\varphi\right)\right\|_{W_{p}^{s, t, q}} \leqslant C_{1} e^{c \alpha}, \tag{2.88}
\end{equation*}
$$

then there exists $C_{2}>0$ such that

$$
\sup _{x \in M}\left|\mathcal{E}_{T, x, \Lambda_{\delta}}(\varphi)\right| \leqslant C_{2}\left(C_{1} T^{\frac{c}{h_{\text {top }}}}+\|\varphi\|_{C^{0}}+1\right) .
$$

If, in addition, Condition 2.4.11 holds for $\delta$, then $\Sigma_{\delta}$ is finite and, taking $\Lambda_{\delta}=\Sigma_{\delta}$ and assuming

$$
t-r+2 \leqslant 0<r-2,
$$

it holds $c=\delta+\epsilon$ in (2.88) for all $\epsilon>0$ and all $\varphi \in W_{p}^{s, t, q}$.
Recall that if $\mathcal{D}_{(\lambda, i, j)} \in W_{p}^{s, t, q}$, for some $\lambda$ with $\Re \lambda>\delta$, is a generalized eigenvector of the generator $X+V$ then for all $\Re \tilde{\lambda}>\delta$ we have that $\mathcal{O}_{\tilde{\lambda}, \tilde{i}, \tilde{j}}\left(\mathcal{D}_{\lambda, i, j}\right)=1$ if $\lambda=\tilde{\lambda}, i=\tilde{\imath}$, and $j=\tilde{\jmath}$, while $\mathcal{O}_{\tilde{\lambda}, \tilde{\imath}, \tilde{j}}\left(\mathcal{D}_{\lambda, i, j}\right)$ vanishes otherwise.

Remark 2.5.8. The condition $\lambda_{\text {max }}=h_{\text {top }}$ is superficial although we show only $\lambda_{\max }=h_{\text {top }}$ and unique simplicity under an additional vanishing assumption in Section 2.5.3. The proof of Theorem 2.5.7 however shows that the horocycle expansion sees only the part of the spectrum with real part below $h_{\text {top }}$ and the eigendistribution $\mu$ which is associated to $h_{\text {top }}$.
Recalling Remark 2.4.8, we find always $\lambda_{\text {min }}<h_{\text {top }}$ if $-s$ and $t$ can be taken to be $1-\epsilon$ for all $\epsilon>0$. This is the case if the geodesic flow is $C^{3-\epsilon}$ for all $\epsilon>0$ (e.g. the flow is of Zygmund type). If one knows then that the weight is $C^{2-\epsilon}$ the basic assumptions of Theorem 2.5.7 are all satisfied (an example is given in Proposition 2.5.10 below for $C^{3}$ contact Anosov flows when $d=3$ ).
Note that $\gamma_{x}\left(\mathcal{D}_{\left(h_{\text {top }, 1,1)}, T\right)}\right.$ is well-defined in the sense of distributions is part of the theorem. By unique ergodicity the expected principal term $T \mu(\varphi)$ is hidden by the term $\gamma_{x}\left(\mathcal{D}_{\left(h_{\text {top }, 1,1)}, T\right)}\right) \mu(\varphi)$ as we ordered the expansion by the distributions $\mathcal{O} \ldots(\varphi)$. We can always write

$$
T \mu(\varphi)=\gamma_{x}(1, T) \mu(\varphi)
$$

and use the expansion result on $\gamma_{x}(1, T)$ again which shows that the leading order term is indeed what we expect. The other terms are modified by the contributions of $\mathcal{O} \ldots(1) \mu(\varphi)$. We make use of this in the following corollary.

Assuming all conditions in the above theorem, this gives polynomial convergence of horocycle averages to the ergodic mean:

Corollary 2.5.9 (Polynomial convergence). Under the assumptions of Theorem 2.5.7 (including Condition 2.4.11 for $\delta$ and $t-r+2 \leqslant 0<r-2$ ) then there exists $\epsilon>0$ such that for all $\varphi \in C_{X}^{r-1}(M)$ there exists $C>0$ such that

$$
\left|\frac{1}{T} \int_{0}^{T} \varphi \circ h_{\rho}(x) \mathrm{d} \rho-\mu(\varphi)\right| \leqslant C T^{-\epsilon},
$$

where $\mu$ is the unique Borel measure which is invariant by the horocycle flow $h_{\rho}$.
Proof. We apply Theorem 2.5.7, using the assumption that Condition 2.4.11 holds for $\delta$ and that $t-r+2 \leqslant 0<r-2$. Then there are only finitely many eigenvalues $\lambda \in \sigma(X+V)$ such that $\Re \lambda>\delta$ and the remainder term $\mathcal{E}_{T, x, \Lambda_{\delta}}(\varphi)$ is bounded from above by $T^{\frac{\delta}{h_{\text {top }}}+\epsilon}$ for all $\epsilon>0$. Hence all but one term in the expansion of the ergodic average decay like $T^{-\epsilon}$ for some $\epsilon>0$. We finally bound the leading term in the expansion

$$
\frac{1}{T} \gamma_{x}\left(\mathcal{D}_{\left(h_{\mathrm{top},}, 1,1\right)}, T\right) \mu(\varphi)-\mu(\varphi)=\frac{1}{T} \gamma_{x}\left(\mathcal{D}_{\left(h_{\mathrm{top}}, 1,1\right)}-1, T\right) \mu(\varphi)
$$

using again Theorem 2.5.7 as before, noting that $\mu(1)=\mu\left(\mathcal{D}_{h_{\text {top }}, 1,1}\right)=1$.
We next discuss the assumptions of our main theorem and the corollary above. We first give sufficient conditions ensuring that $\partial_{\rho} \tau(0,-\alpha, \cdot) \in C^{r-1}$ and that there exist parameters in our anisotropic space giving $\lambda_{\min }(s, t, p)<h_{\text {top }}$ :

Proposition 2.5.10. Let $g_{\alpha}$ be a $C^{3}$ contact Anosov flow on a closed Riemannian manifold $M$ of dimension $d=3$ preserving a $C^{1}$ contact form and let the strong-stable distribution $E_{-}$be orientable. Then there exists a horocycle flow $h_{\rho}$ such that $\partial_{\rho} \tau(0,-\alpha, \cdot) \in C^{r-1}$ for every $\alpha \geqslant 0$ and for any $r \in[2,3)$.
Setting $-s=t=\frac{r-1}{2}-\frac{\epsilon}{2}$ for suitable $0<\epsilon<\frac{r-1}{2}$, the constant $\lambda_{\text {min }}(s, t, p)$ is independent of $p$ and can be taken arbitrary close to $0^{+}$while $t-r+2 \leqslant 0<r-2$.

Proof. The contact assumption means that there is an invariant 1-form $\eta \in T^{*} M$ such that $\mu:=\eta \wedge \mathrm{d} \eta \neq 0$ everywhere. By assumption $\eta$ is $C^{1}$. Moreover $\eta$ is annihilated on $E_{+}+E_{-}$and $\mu \in \wedge^{3} T^{*} M$ is preserved by the flow. We use [40, Theorem 3.1] together with the comment on the relation between Zygmund
and Hölder regularity to infer that the strong-stable distribution is $C^{r-1}$ for all $r \in[2,3)$ if $d=3$. Hence for the horocycle flow given by the unit speed parametrization (and more general all of its $C^{r-1}$ reparametrizations) we find $\partial_{\rho} \tau(0,-\alpha, \cdot) \in C^{r-1}$. By assumption the flow $g_{\alpha}$ preserves volume and $d=3$.
To see a gap between $\lambda_{\min }=\lambda_{\min }(s, t, p)$ and $h_{\text {top }}$, we may assume the unit speed parametrization of the horocycle flow $h_{\rho}$. It follows by Proposition 2.5.13 (viii) and Lemma 2.3.3 that for all $C^{r-1}$ reparametrizations the resulting transfer operators are conjugate to each other.
Then it follows from Lemma 2.4.7 together with Proposition 2.5.13 (viii) that $\lambda_{\text {min }}$ is independent of $p$ and is arbitrary close to $0^{+}$for a suitable choice of $s, t$ and $r$. Moreover, if we assume $0<t \leqslant \frac{r-3}{2}+\epsilon$ we satisfy $t-r+2 \leqslant 0<r-2$ since $\epsilon<\frac{r-1}{2}$.

Second, we discuss Condition 2.4.11:
Remark 2.5.11. Condition 2.4.11 was inspired by estimates of Dolgopyat [23], who was working with operators acting on symbolic spaces. This condition, replacing however our $W_{p}^{s, t, q}$ by other anisotropic Banach spaces, was proved by several authors [9], [31], [48], [66] for the generator $X+V$, associated to contact Anosov flows and $V=0$ the trivial potential, for which they also obtained the additional condition in Corollary 2.5.9.

In the case of geodesic flows on compact surfaces of constant negative curvature, we find that $V$ is a constant, but the fact that our Banach space is different makes it difficult to apply the results of [9], [31], [48], [66] directly in order to establish Condition 2.4.11. We expect however that the condition holds and (as pointed out by Liverani and Butterley) can be obtained by exploiting e.g. [19, Remark 2.6].

For non-constant potential $V$, since Dolgopyat [23] obtained exponential decay of correlations for Gibbs measures with Hölder potentials, we expect that Condition 2.4.11 indeed holds also in our setting, in particular for compact surfaces of variable negative curvature (e.g. using an argument similar as for the proof in [24, Proposition 3.4]). (We warn the reader that the value of $\delta$ given by Dolgopyat-type arguments is usually very close to $\lambda_{\text {max }}$.)

We end this subsection by a comparison of our main theorem and the results of Flaminio and Forni [28]: Let $M$ be the unit tangent bundle of a compact hyperbolic Riemann surface. Let $g_{\alpha}$ be its unit speed geodesic flow and let vol be the canonical (invariant) volume form on $M$ (which is also a measure of
maximal entropy) and consider the unit speed horocycle flow which leaves vol invariant as well (hence $\mu=\operatorname{vol}$ ). Then $h_{\text {top }}=1$ because $\tau(\rho, \alpha, x)=\rho \exp (-\alpha)$ and $\mathcal{D}_{(1,1,1)} \equiv 1$ (hence $\left.\gamma_{x}\left(\mathcal{D}_{(1,1,1)}, T\right)=T\right)$. In the setting of Riemann surfaces, the possible Jordan blocks are known [28, Theorem 1.5]. In particular, the eigenvalue $h_{\text {top }}=1$ is simple, there are no other eigenvalues of real part equal to one, all eigenvalues with $\Re \lambda>0$ are semi-simple, and there are only finitely many eigenvalues with $\Re \lambda>\frac{1}{2}$. Moreover, since the vector fields are constant, the regularity parameters $-s, t$ can be taken large enough such that $\lambda_{\min }<0$. Hence we can take any $\delta \geqslant 0$ in Theorem 2.5.7, and we find, for any finite subset of $\Sigma_{\delta}$ containing 1,

$$
\int_{0}^{T} \varphi \circ h_{\rho}(x) \mathrm{d} \rho=T \operatorname{vol}(\varphi)+\sum_{\lambda \in \Lambda_{\delta} \backslash\{1\}} \sum_{i=1}^{n_{\lambda}} T^{\lambda} c_{(\lambda, i, 1)}(T, x) \mathcal{O}_{(\lambda, i, 1)}(\varphi)+\mathcal{E}_{T, x, \Lambda_{\delta}}(\varphi)
$$

where we can take $\Lambda_{\delta}=\Sigma_{\delta}$ if $\delta \geqslant \frac{1}{2}$, and where $c_{(\lambda, i, 1)}$ and $\mathcal{E}_{T, x, \Lambda_{\delta}}$ satisfy the claims of Theorem 2.5.7 (with an additional $\log T$-factor if $\lambda=\frac{1}{2}$ ). In particular, if Condition 2.4.11 holds for some $\delta>\frac{1}{2}$ (see Remark 2.5.11) there exists $C>0$ such that for all $\epsilon>0$

$$
\left|\mathcal{E}_{T, x, \Sigma_{\delta}}(\varphi)\right| \leqslant C T^{\delta+\epsilon} .
$$

Note that we required $\delta \geqslant 0$ because for $\delta<0$ we find no improvement of the remainder term (this comes the local bounds in Lemma 2.5.14). An analogous behavior is seen in the corresponding expansion of Flaminio-Forni in [28, Theorem 1.5]. However they are not limited to finite sets $\Lambda_{\delta}$ of eigenvalues (Faure-Tsujii do not seem to be limited either in [26]). Our methods, however, do not seem to allow to go beyond the first vertical line with infinitely many resonances in $\left.\sigma(X+V)\right|_{W_{p}^{s, t, q}}$ in the expansion of the horocycle integral. (This could be a natural limitation, as discussed in [66, p.1497, below Theorem 1.1].)

### 2.5.2 Weighted horocycle integrals, properties of $\tau$, local bounds

In order to use a smooth cutoff trick of Giulietti-Liverani to decompose $\gamma_{x}(\cdot, T)$ in Lemma 2.5.14 below, we need to consider weighted horocycle integrals: For all $\varphi \in C_{X}^{r-1}(M)$, for all compactly supported $w \in C(\mathbb{R}, \mathbb{C})$ and for all $x \in M$, let

$$
\begin{equation*}
\gamma_{w, x}(\varphi):=\int_{\mathbb{R}} w(\rho) \cdot\left(\varphi \circ h_{\rho}(x)\right) \mathrm{d} \rho . \tag{2.89}
\end{equation*}
$$

denote the horocycle integral of the horocycle flow $h_{\rho}$ for the observable $\varphi$ at base point $x$ with weight $w$.
For further purposes, it is useful to view $\gamma_{w, x}$ as a functional in the topological dual space of $W_{p}^{s, t, q}$ for weights $w$ with compact support and sufficient differentiability:

Lemma 2.5.12. Let $p \in[1, \infty]$ and let $0<q \leqslant t<r-1$ and let $-r<s<0$. Let $x \in M$. Then for some $C>0$, for all $C^{-s}$ maps $w: \mathbb{R} \rightarrow \mathbb{C}$ with compact support it holds

$$
\left\|\gamma_{w, x}\right\|_{W_{p}^{s, t, q} \rightarrow \mathbb{C}} \leqslant C|\operatorname{supp} w|\|w\|_{C^{-s}} .
$$

Proof. We recall the partition of unity $\vartheta_{\varpi}$ and chart maps $\kappa_{\varpi}, \varpi \in \Omega$ (see Definition 2.3.6). We set for all $x \in M$, for all $\alpha \geqslant 0$ and for all $\varphi \in C_{X}^{r-1}(M)$

$$
\begin{align*}
y_{x, \varpi_{1}, \alpha}(\rho) & :=\kappa_{\varpi_{1}} \circ g_{\alpha} \circ h_{\rho}(x),  \tag{2.90}\\
\varphi_{w, x, \varpi_{1}, \varpi_{2}, \alpha}(z) & :=\left(\vartheta_{\varpi_{2}} \cdot \phi_{-\alpha} \circ g_{-\alpha}\right) \circ \kappa_{\varpi_{1}}^{-1}(z) \cdot \int_{-\infty}^{\infty} w(\rho) \delta\left(z-y_{x, \varpi, \alpha}(\rho)\right) \mathrm{d} \rho, \\
\varphi_{\varpi_{1}, \alpha}(z) & :=\left(\vartheta_{\varpi_{1}} \cdot \mathcal{L}_{\alpha, \phi_{\alpha}} \varphi_{1}\right) \circ \kappa_{\varpi_{1}}^{-1}(z) .
\end{align*}
$$

With this notation, recalling the weighted horocycle integral associated to Definition 2.5.1, we express for all $\alpha \geqslant 0$

$$
\begin{equation*}
\gamma_{w, x}(\varphi)=\sum_{\varpi_{1}, \varpi_{2} \in \Omega} \int_{\mathbb{R}^{d}} \varphi_{w, x, \varpi_{1}, \varpi_{2}, \alpha}(z) \cdot \varphi_{\varpi_{1}, \alpha}(z) \mathrm{d} z . \tag{2.91}
\end{equation*}
$$

We set

$$
c^{\prime}(+):=-s, \quad c^{\prime}(0):=-t, \quad c^{\prime}(-):=-t .
$$

We recall $\widetilde{\Psi}_{\sigma, n}$ defined in (2.24). We bound, using Plancherel's Theorem, CauchySchwarz for the sum in $\sigma$ and $n$, and twice Hölder's inequality with respect to $z$ and $\alpha$, respectively, for some constant $C>0$

$$
\begin{align*}
\alpha_{0}\left|\gamma_{w, x}(\varphi)\right| & =\int_{0}^{\alpha_{0}}\left|\sum_{\varpi_{1}, \varpi_{2}} \int_{\mathbb{R}^{d}} \varphi_{w, x, \varpi_{1}, \varpi_{2}, \alpha}(z) \cdot \varphi_{\varpi_{1}, \alpha}(z) \mathrm{d} z\right| \mathrm{d} \alpha \\
\leqslant & \int_{0}^{\alpha_{0}} \sum_{\varpi_{1}, w_{2}}\left|\int_{\mathbb{R}^{d}} \sum_{\sigma, n} 2^{-c(\sigma) n} \widetilde{\Psi}_{\sigma, n}^{\mathrm{Op}} \varphi_{w, x, \varpi_{1}, \varpi_{2}, \alpha}(z) 2^{c(\sigma) n} \Psi_{\sigma, n}^{\mathrm{Op}} \varphi_{\varpi_{1}, \alpha}(z) \mathrm{d} z\right| \mathrm{d} \alpha \\
& \leqslant C \sup _{\alpha, \varpi_{1}, \varpi_{2}}\left\|\left(\sum_{\sigma, n} 4^{c^{\prime}(\sigma) n}\left|\widetilde{\Psi}_{\sigma, n}^{\mathrm{Op}} \varphi_{w, x, \varpi_{1}, \varpi_{2}, \alpha}\right|^{2}\right)^{\frac{1}{2}}\right\|_{L_{p^{*}}}\|\varphi\|_{W_{p}^{s, t, q}}, \tag{2.92}
\end{align*}
$$

where $p^{*}:=1-\frac{1}{p}$ is the Hölder conjugate of $p$. To conclude, it is enough to establish the following upper bounds for $\left\|\widetilde{\Psi}_{\sigma, n}^{\mathrm{Op}} \varphi_{w, x, \omega_{1}, \omega_{2}, \alpha}\right\|_{L_{p^{*}}}$ :
(i) There exists a constant $C_{1}>0$ such that for every $C^{0}$ map $w: \mathbb{R} \rightarrow \mathbb{R}$, every $p \in[1, \infty]$, every $\sigma \in\{-,+, 0\}, n \in \mathbb{N}$, every $\varpi_{1}, \varpi_{2} \in \Omega$ it holds

$$
\left\|\widetilde{\Psi}_{\sigma, n}^{\mathrm{Op}_{n}} \varphi_{w, x, \varpi_{1}, \varpi_{2}, \alpha}\right\|_{L_{p}} \leqslant C_{1}|\operatorname{supp} w|\|w\|_{L_{\infty}}, \forall x \in M, \forall 0 \leqslant \alpha \leqslant \alpha_{0} .
$$

(ii) There exists a constant $C_{2}>0$ such that for every $-r<s<0$, for every $C^{|s|}$ map $w: \mathbb{R} \rightarrow \mathbb{R}$ with compact support, every $p \in[1, \infty]$, every $n \in \mathbb{N}$, every $\varpi_{1}, \varpi_{2} \in \Omega$ it holds

$$
\left\|\widetilde{\Psi}_{-, n}^{\mathrm{Op}} \varphi_{w, x, w_{1}, w_{2}, \alpha}\right\|_{L_{p}} \leqslant C_{2} 2^{s n}|\operatorname{supp} w|\|w\|_{C^{-s}}, \quad \forall x \in M, \forall 0 \leqslant \alpha \leqslant \alpha_{0} .
$$

We first show claim (i). We fix $w, \sigma, n, \varpi_{1}, \varpi_{2}, x$ and $\alpha$. We let $J \subseteq \operatorname{supp} w$ be the maximal subset such that $y_{x, w_{1}, \alpha \mid J}$ is well-defined. We note that $J$ decomposes into a finite disjoint union, e.g. $J=\bigsqcup_{k=1}^{N} I_{k}$ for some $N \in \mathbb{N}$ and some real intervals $I_{k}$. In particular, since the flow $h_{\rho}$ is non-singular and, in addition the manifold $M$ is compact and each stable leaf is dense in $M$ and $0 \leqslant \alpha \leqslant \alpha_{0}$, for some constant $C_{1}>0$, we have $\left|I_{k}\right| \leqslant C_{1} \operatorname{diam} V_{\varpi_{1}}$ and $N \leqslant C_{1} \frac{\mid \text { supp } w \mid}{\text { diam } V_{\varpi_{1}}}$. For every $z \in \mathbb{R}^{d}$ we estimate for some constants $C_{2}, \ldots, C_{4}>0$

$$
\begin{align*}
& \left|\varphi_{w, x, \varpi_{1}, \varpi_{2}, \alpha}(z)\right| \leqslant C_{2}\left|\int_{J} w(\rho) \delta\left(z-y_{x, \varpi_{1}, \alpha}(\rho)\right) \mathrm{d} \rho\right| \\
& =\left|\sum_{k=1}^{N} \int_{I_{k}} w(\rho) \delta\left(z-y_{x, \varpi_{1}, \alpha}(\rho)\right) \mathrm{d} \rho\right|=\left.\left|\sum_{k=1}^{N} \sum_{\rho \in y^{-1}(z) \cap I_{k}} w(\rho)\right| \partial_{\rho} y_{x, \varpi_{1}, \alpha}(\rho)\right|^{-1} \mid \\
& \leqslant\left. C_{3} N \max _{\rho \in \operatorname{supp} w}|w(\rho)| \partial_{\rho} y_{x, \varpi_{1}, \alpha}(\rho)\right|^{-1}\left|\leqslant C_{4}\right| \operatorname{supp} w \mid\|w\|_{L_{\infty}}, \tag{2.93}
\end{align*}
$$

where we used in the last step non-singularity of $h_{\rho}$ and $0 \leqslant \alpha \leqslant \alpha_{0}$. We conclude, using Young's inequality on $\left\|\widetilde{\Psi}_{\sigma, n}^{\mathrm{Op}} \varphi_{w, x, \varpi_{1}, \varpi_{2}, \alpha}\right\|_{L_{p}}$ together with the bound in (2.93).
We now show claim (ii). Again we fix $w, \sigma, n, \varpi_{1}, \varpi_{2}, x$ and $\alpha$ and set $y:=$ $y_{x, \varpi_{1}, \alpha}$. Analogously as in the proof of (i), we let $I_{k} \subseteq \mathbb{R}, 1 \leqslant k \leqslant N$, be the $N$ connected components of $J$ for some $N \in \mathbb{N}$. For every $z \in \mathbb{R}^{d}$ we expand

$$
\widetilde{\Psi}_{-, n}^{\mathrm{Op}} \varphi_{w, x, \varpi_{1}, \varpi_{2}, \alpha}(z)=\frac{2^{d n}}{(2 \pi)^{d}} \int_{\mathbb{R}^{d}} \int_{\mathbb{R}^{d}} \tilde{\Psi}_{-, 1}(\xi) e^{\mathrm{i} 2^{n} \xi(z-\widetilde{y})} \varphi(\widetilde{y}) \mathrm{d} \xi \mathrm{~d} \tilde{y},
$$

where we set

$$
\varphi(\widetilde{y}):=\varphi_{w, x, w_{1}, w_{2}, \alpha}(\widetilde{y}) .
$$

We note that $\operatorname{supp} \varphi \subseteq y(J)$. In particular, we reparametrize $\tilde{y} \in \operatorname{supp} \varphi$ by $\widetilde{y}=\tilde{z}(\widetilde{\rho})$ for some diffeomorphism $z \in C^{r}$ and $\tilde{\rho} \in \mathbb{R}$. We set $\mathrm{D}_{\tilde{\rho}}(\cdot):=\mathrm{i} \partial_{\tilde{\rho}} \frac{(\cdot)}{\partial_{\tilde{\rho}} \tilde{z}}$. Since $\tilde{z}(\mathbb{R})$ is a piece of a stable manifold in charts there exists a constant $C_{3}>0$ such that we have $\left|\partial_{\tilde{\rho}} \xi \widetilde{z}(\widetilde{\rho})\right| \geqslant C_{1} 2^{n}$ for all $\xi$ in supp $\widetilde{\Psi}_{-, n}$ is essentially part of an unstable cone in charts by construction. We note that $\varphi \circ z$ is $C^{r}$. Using $[-s]-$ times integration by parts (see Lemma B.3), followed by a regularized integration by parts with respect to $\widetilde{\rho}$ if $-s \notin \mathbb{N}$, respectively (see Lemma B. 5 in which we take $d=1, G=y$ and $L^{-1}=\epsilon=2^{-n}$ ), this yields

$$
\widetilde{\Psi}_{-, n}^{\mathrm{Op}} \varphi_{w, x, \varpi_{1}, \varpi_{2}, \alpha}(z)=\frac{2^{((s+d) n)}}{(2 \pi)^{d}} \int_{\mathbb{R}^{d}} \int_{\mathbb{R}} \widetilde{\Psi}_{+, 1}(\xi) e^{\mathrm{i} 2^{n} \xi(z-\widetilde{z}(\widetilde{\rho}))} \widetilde{\mathrm{D}}_{\widetilde{\rho}}^{-s} \widetilde{\varphi}(\widetilde{\rho}) \mathrm{d} \xi \mathrm{~d} \widetilde{\rho},
$$

where

$$
\begin{aligned}
\widetilde{\mathrm{D}}_{\widetilde{\rho}}^{-s} \tilde{\varphi}: & :\left\{\begin{array}{ll}
\mathrm{D}_{\widetilde{\rho}}^{-s}\left(\tilde{\varphi} \circ \tilde{z} \cdot \partial_{\widetilde{\rho}} \tilde{z}\right), & \text { if }-s \in \mathbb{N} \\
2^{(-s+\lfloor-s\rfloor) n}\left(\partial_{\tilde{\rho}}\left(\frac{1}{2^{n} \tilde{z}} \widetilde{\widetilde{\varphi}}\right)_{\epsilon}+\partial_{\tilde{\rho}} \tilde{z} \cdot\left(\frac{1}{\xi \widetilde{\xi}} \widetilde{\tilde{\varphi}}-\left(\frac{1}{\xi \tilde{z}} \widetilde{\tilde{\varphi}}\right)_{\epsilon}\right)\right), & \text { if }-s \notin \mathbb{N}
\end{array},\right. \\
\widetilde{\varphi}: & =\mathrm{D}_{\widetilde{\rho}}^{\lfloor-s\rfloor}\left(\tilde{\varphi} \circ \tilde{z} \cdot \partial_{\widetilde{\rho}} \widetilde{z}\right),
\end{aligned}
$$

and the $\epsilon$-term is just the convolution $\left(\frac{1}{\xi y} \mathrm{D}_{\tilde{\rho}}^{[-s]} \tilde{\varphi}\right) * \nu_{\epsilon}$ with a $C^{\infty}$ map $\nu_{\epsilon}$ with $\operatorname{supp} \nu_{\epsilon} \subseteq(-\epsilon, \epsilon)$. Note that all derivatives of $\tilde{\mathrm{D}}_{\widetilde{\rho}}^{-s} \tilde{\varphi}$ in $\xi$ are bounded in $n$, using Lemma B. 5 and non-singularity of $h_{\rho}$ and $0 \leqslant \alpha \leqslant \alpha_{0}$. We proceed analogously as in the proof of Lemma 2.3.3, integrating $(d+1)$-times by parts in $\xi$ if $2^{n}|z-\widetilde{z}(\widetilde{\rho})|>1$ and conclude, using that $\operatorname{supp} \tilde{\mathrm{D}}_{\widetilde{\rho}}^{-s} \widetilde{\varphi} \subseteq \operatorname{supp} \widetilde{\varphi}$ is bounded.

We group below some properties of the pointwise renormalization. (Note in particular that Claim (xi) in Proposition 2.5.13, which will follow from [31, Remark C.4] of Giulietti-Liverani-Pollicott, will play a key part to estimate the spectral bound of $X+V$. Also, Claim (viii) in Proposition 2.5.13 shows that $\phi_{\alpha}=\partial_{\rho} \tau(0,-\alpha, \cdot)$ differs from the unit speed parametrization function by a multiplicative 1-coboundary.)

Proposition 2.5.13 (Properties of pointwise renormalization). Let $\tau$ be the renormalization time of a stable horocycle flow. For all $\rho, \alpha \in \mathbb{R}$ and for all $x \in M$ it holds:
(i) $\tau(0, \alpha, x)=0$,
(ii) $\tau(\rho, 0, x)=\rho$,
(iii) $\tau\left(\rho, \alpha_{1}+\alpha_{2}, x\right)=\tau\left(\tau\left(\rho, \alpha_{2}, x\right), \alpha_{1}, g_{\alpha_{2}}(x)\right)$, for all $\alpha_{1}, \alpha_{2} \in \mathbb{R}$,
(iv) $\tau\left(\rho_{1}+\rho_{2}, \alpha, x\right)=\tau\left(\rho_{1}, \alpha, h_{\rho_{2}}(x)\right)+\tau\left(\rho_{2}, \alpha, x\right)$, for all $\rho_{1}, \rho_{2} \in \mathbb{R}$,
(v) $\partial_{\rho} \tau(\rho, \alpha, x)=\partial_{\rho} \tau\left(0, \alpha, h_{\rho}(x)\right)$,
(vi) $\partial_{\rho} \tau\left(0, \alpha_{1}, g_{\alpha_{2}}(x)\right) \partial_{\rho} \tau\left(0, \alpha_{2}, x\right)=\partial_{\rho} \tau\left(0, \alpha_{1}+\alpha_{2}, x\right)$, for all $\alpha_{1}, \alpha_{2} \in \mathbb{R}$,
(vii) $\tau(\rho, \alpha, x)=\gamma_{x}\left(\partial_{\rho} \tau(0, \alpha, \cdot), \rho\right)$,
(viii) $\partial_{\rho} \tau(0, \alpha, \cdot)=\operatorname{det} \mathrm{D} g_{\alpha \mid E_{-}} \frac{\left(\partial_{\rho} h_{0}\right)^{*}\left(\partial_{\rho} h_{0}\right)}{\left(\partial_{\rho} h_{0} g_{\alpha}\right)^{*}\left(\partial_{\rho} h_{0} 0 g_{\alpha}\right)}$ where $\partial_{\rho} h_{0}(x):=\partial_{\rho} h_{\rho \mid \rho=0}(x)$,
(ix) $0<\partial_{\rho} \tau(0, \alpha, x)<\infty$,
(x) if $\alpha \geqslant 0$ there exist $C_{1}>0$ and $0<\theta<1$ both independent of $\alpha$ and $x$ such that $\left\|\partial_{\rho} \tau(\cdot, \alpha, x)\right\|_{C^{r-1}} \leqslant C_{1} \theta^{\alpha}$,
(xi) if $|\rho| \geqslant 1$ and $\alpha \geqslant 0$ there exists $C_{2} \geqslant 1$ independent of $\rho, \alpha$ and $x$ such that

$$
C_{2}^{-1}|\rho| e^{h_{\text {top }} \alpha} \leqslant|\tau(\rho,-\alpha, x)| \leqslant C_{2}|\rho| e^{h_{\text {top }} \alpha} .
$$

(xii) if $\alpha \geqslant 0$ such that $|\tau(\rho, \alpha, x)|=c$ for some $c \geqslant 1$ then there exists $C_{3} \geqslant 1$ independent of $\rho, \alpha$ and $x$ such that

$$
C_{3}^{-1} c e^{h_{\text {top } \alpha}} \leqslant|\rho| \leqslant C_{3} c e^{h_{\text {top }} \alpha} .
$$

Proof. We note that by Definition 2.5.1, Definition 2.5.4 and Lemma 2.5.6 the renormalization time is differentiable in $\rho$. Every stable leaf is dense in $M$ hence together with non-singularity of the flow $h_{\rho}$ it follows $h_{\rho_{1}}(x)=h_{\rho_{2}}(x) \Rightarrow \rho_{1}=\rho_{2}$. Then Claim (i)-(ii) follow directly from (2.84). We deduce from (2.84)

$$
\begin{aligned}
h_{\tau\left(\rho, \alpha_{1}+\alpha_{2}, x\right)} \circ g_{\alpha_{1}+\alpha_{2}}(x) & =g_{\alpha_{1}+\alpha_{2}} \circ h_{\rho}(x)=g_{\alpha_{1}} \circ h_{\tau\left(\rho, \alpha_{2}, x\right)} \circ g_{\alpha_{2}}(x) \\
& =h_{\tau\left(\tau\left(\rho, \alpha_{2}, x\right), \alpha_{1}, g_{\alpha_{2}}(x)\right)} \circ g_{\alpha_{1}+\alpha_{2}}(x)
\end{aligned}
$$

This yields Claim (iii). Also from (2.84) we find

$$
\begin{aligned}
h_{\tau\left(\rho_{1}+\rho_{2}, \alpha, x\right)} \circ g_{\alpha}(x) & =g_{\alpha} \circ h_{\rho_{1}+\rho_{2}}(x)=g_{\alpha} \circ h_{\rho_{1}} \circ g_{-\alpha} \circ g_{\alpha} \circ h_{\rho_{2}} \circ g_{-\alpha} \circ g_{\alpha}(x) \\
& =h_{\tau\left(\rho_{1}, \alpha, h_{\rho_{2}}(x)\right)} \circ h_{\tau\left(\rho_{2}, \alpha, x\right)} \circ g_{\alpha}(x) .
\end{aligned}
$$

This yields Claim (iv). Claim (v) and (vi), using Claim (i), follow by differentiating both sides in (iv) and (iii) at $\rho_{1}=0$ and $\rho=0$, respectively.
Claim (vii) follows from (2.83) and (v).
To show Claim (viii), we take derivatives on both sides of (2.84) with respect to $\rho$

$$
\begin{equation*}
\mathrm{D} g_{\alpha} \partial_{\rho} h_{\rho}(x)=\partial_{\rho} \tau(\rho, \alpha, x) \cdot\left(\partial_{\rho} h_{0}\right) \circ h_{\tau(\rho, \alpha, x)} \circ g_{\alpha}(x) . \tag{2.94}
\end{equation*}
$$

2 Horocycle averages on closed manifolds

Now we let $\left(\partial_{\rho} h_{0}\right)^{*} \in E_{-}^{*}$ be the canonical dual of $\partial_{\rho} h_{0}$. We calculate

$$
\begin{gather*}
\left(\partial_{\rho} h_{0} \circ g_{\alpha}\right)^{*}\left(\mathrm{D} g_{\alpha} \partial_{\rho} h_{0}\right)=\left(\partial_{\rho} h_{0} \circ g_{\alpha}\right)^{*}\left(\left(g_{\alpha}\right)_{*} \partial_{\rho} h_{0}\right)=\left(g_{\alpha}\right)^{*}\left(\partial_{\rho} h_{0} \circ g_{\alpha}\right)^{*}\left(\partial_{\rho} h_{0}\right) \\
=\operatorname{det}\left(\mathrm{D} g_{\alpha \mid E_{-}}\right)^{*}\left(\partial_{\rho} h_{0}\right)^{*}\left(\partial_{\rho} h_{0}\right)=\operatorname{det} \mathrm{D} g_{\alpha \mid E_{-}}\left(\partial_{\rho} h_{0}\right)^{*}\left(\partial_{\rho} h_{0}\right) . \tag{2.95}
\end{gather*}
$$

We set $\rho=0$ in (2.94) and conclude, using (2.95) and non-singularity of the horocycle flow.
Claim (ix) follows from (viii) together with the fact $\lim _{\alpha \rightarrow 0} \operatorname{det} \mathrm{D} g_{\alpha \mid E-}=1$ and compactness of $M$.
In order to show (x), we note first, since $r \geqslant 2$, using Claim (v) and the cocycle property (vi),

$$
\partial_{\rho} \tau(\rho, \alpha, x)=\partial_{\rho} \tau\left(0, \alpha, h_{\rho}(x)\right)=\exp -\int_{0}^{\alpha} V \circ g_{\widetilde{\alpha}} \circ h_{\rho}(x) \mathrm{d} \tilde{\alpha},
$$

where $V:=-\partial_{\alpha} \partial_{\rho} \tau(0,0, \cdot) \in C^{r-1}$. Therefore it holds, using the equality in (2.84),

$$
\begin{align*}
\partial_{\rho}^{2} \tau(\rho, \alpha, x) & =-\partial_{\rho} \tau(\rho, \alpha, x) \cdot \partial_{\rho} \int_{0}^{\alpha} V \circ g_{\tilde{\alpha}} \circ h_{\rho}(x) \mathrm{d} \tilde{\alpha} \\
& =-\partial_{\rho} \tau(\rho, \alpha, x) \cdot \int_{0}^{\alpha} \partial_{\rho} \tau(\rho, \widetilde{\alpha}, x) \cdot\left(\mathrm{D} V \partial_{\rho} h_{0}\right) \circ g_{\widetilde{\alpha}} \circ h_{\rho}(x) \mathrm{d} \tilde{\alpha}, \tag{2.96}
\end{align*}
$$

where $\left|\partial_{\rho} \tau(\rho, \alpha, x)\right| \leqslant C \theta^{\alpha}$ for some $0<\theta<1$ and $C_{1}>0$ both independent of $\alpha, \rho$ and $x$ by (viii). Hence there is $C_{2}=C_{2}(V)>0$ such that $\left|\partial_{\rho}^{2} \tau(\rho, \alpha, x)\right| \leqslant$ $C_{2} \theta^{\alpha}$. By induction, all derivatives $\partial_{\rho}^{k} \tau(\rho, \alpha, x)$, where $k \in \mathbb{N}$, depend only on $\partial_{\rho} \tau(\rho, \alpha, x)$ (and $k$ and derivatives of $V$ which are independent of $\alpha$ ) and so does the Hoelder norm $\left\|\partial_{\rho} \tau(\rho, \alpha, x)\right\|_{C^{r-1}}$. Since $r \geqslant 2$ the Hoelder coefficicent of $\partial_{\rho} \tau(\cdot, \alpha, x)$ is bounded by $\left\|\partial_{\rho}^{2} \tau(\cdot, \alpha, x)\right\|_{C^{0}}$ and we conclude.
Claim (xi) for $\rho \geqslant 1$ and $\alpha \leqslant 0$ follows from [31, Lemma C.1] and [31, Remark C.4] (recall that $g_{\alpha}$ is transitive) in which we replace $W$ with a manifold which contracts in forward time. To this end we set $W_{x}:=h_{[0,1]}(x)$ for every $x \in M$. Since the stable flow is non-singular, the stable manifold $W_{x}$ is of bounded length (from above and below) for all $x \in M$. We estimate, using Proposition 2.5.13 (viii) for the first and [31, Remark C.4] for the last inequality, with constants
$C_{3}, \ldots, C_{6}>0$ independent of $\rho, \alpha, x$

$$
\begin{aligned}
\tau(\rho,-\alpha, x) & \leqslant C_{3} \int_{0}^{\rho} \operatorname{det} \mathrm{D} g_{-\alpha \mid E_{-}} \circ h_{\rho}(x) \mathrm{d} \rho \leqslant C_{4} \rho \int \operatorname{det} \mathrm{D} g_{-\alpha \mid E_{-}} \mathrm{d} W_{x} \\
& \leqslant C_{5} \rho \operatorname{vol}\left(g_{-\alpha}\left(W_{x}\right)\right) \leqslant C_{6} \rho e^{h_{\text {top }} \alpha} .
\end{aligned}
$$

A lower bound for $\tau(\rho,-\alpha, x)$ is obtained in an analogous way, using the last statement in [31, Lemma C.1]. We conclude for all $|\rho| \geqslant 1$, noting that $\tau(-\rho, \alpha, x)=$ $-\tau\left(\rho, \alpha, h_{-\rho}(x)\right)$, using Claims (iv) and (i).
Claim (xii) follows from Proposition 2.5.13 (viii), and the following equality which follows from Claim (iii)

$$
\rho=\tau\left(\tau(\rho, \alpha, x),-\alpha, g_{\alpha}(x)\right)=\tau\left(c,-\alpha, g_{\alpha}(x)\right) .
$$

We shall use in the next two lemmas the following key identity for the horocyle integral (2.89)

$$
\begin{equation*}
\gamma_{w, x}(\varphi)=\gamma_{w \circ \tau\left(,,-\alpha, g_{\alpha}(x)\right), g_{\alpha}(x)}\left(\mathcal{L}_{\alpha, \partial_{\rho} \tau(0,-\alpha, \cdot)}\right), \forall \alpha \geqslant 0 . \tag{2.97}
\end{equation*}
$$

To check the above identity, using (2.84) and Proposition 2.5.13 (iii), (v)-(vi), just notice that for all $\alpha \in \mathbb{R}$

$$
\begin{align*}
\gamma_{w, x}(\varphi) & =\int_{-\infty}^{\infty} w(\rho) \cdot \varphi \circ g_{-\alpha} \circ g_{\alpha} \circ h_{\rho}(x) \mathrm{d} \rho \\
& =\int_{-\infty}^{\infty} w(\rho) \cdot \varphi \circ g_{-\alpha} \circ h_{\tau(\rho, \alpha, x)} \circ g_{\alpha}(x) \mathrm{d} \rho \\
& =\int_{-\infty}^{\infty} w\left(\tau\left(\rho,-\alpha, g_{\alpha}(x)\right)\right) \cdot \varphi \circ g_{-\alpha} \circ h_{\rho} \circ g_{\alpha}(x) \cdot \partial_{\rho} \tau\left(\rho,-\alpha, g_{\alpha}(x)\right) \mathrm{d} \rho \\
& =\int_{-\infty}^{\infty} w\left(\tau\left(\rho,-\alpha, g_{\alpha}(x)\right)\right) \cdot\left(\partial_{\rho} \tau(0,-\alpha, \cdot) \cdot \varphi \circ g_{-\alpha}\right) \circ h_{\rho} \circ g_{\alpha}(x) \mathrm{d} \rho \\
& =\gamma_{w \circ \tau\left(\cdot,-\alpha, g_{\alpha}(x)\right), g_{\alpha}(x)}\left(\partial_{\rho} \tau(0,-\alpha, \cdot) \cdot \varphi \circ g_{-\alpha}\right)  \tag{2.98}\\
& =\gamma_{w \circ \tau\left(\cdot,-\alpha, g_{\alpha}(x)\right), g_{\alpha}(x)}\left(\mathcal{L}_{\left.\alpha, \partial_{\rho} \tau(0,-\alpha, \cdot) \varphi\right), \quad \text { if } \alpha \geqslant 0 .}\right.
\end{align*}
$$

We now state upper bounds for $\left|\gamma_{x}(\varphi, T)\right|$ similar to the results in [28, Lemma 5.16]. The prove uses the analogue of the smooth cutoff used by GiuliettiLiverani [30] but uses a different construction of the local decomposition of $\gamma_{x}(\varphi, T)$.
Lemma 2.5.14 (Local bounds). For every $T>0$ and for every $x \in M$ there exists $w \in C^{-s}(\mathbb{R},[0,1])$ such that for every $\varphi \in W_{p}^{s, t, q}$, where $p \in[1, \infty]$ and

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$s<0<q \leqslant t<r-1+s$, the following holds:
(i) There exists $C_{1}>0$ independent of $T, x$ and $\varphi$ such that

$$
\left|\gamma_{x}(\varphi, T)-\gamma_{w, x}(\varphi)\right| \leqslant C_{1}\|\varphi\|_{C^{0}} .
$$

Moreover, if $-s<\frac{\theta_{\text {min }}}{\theta_{\text {max }}}$, where for some $0<\theta_{\min }<\theta_{\max }$, for some $C_{0} \geqslant 1$ and for all $\alpha \geqslant 0$ it holds

$$
\begin{equation*}
C_{0}^{-1} e^{-\theta_{\max } \alpha} \leqslant \inf _{x \in M} \partial_{\rho} \tau(\cdot, \alpha, x) \leqslant \sup _{x \in M} \partial_{\rho} \tau(\cdot, \alpha, x) \leqslant C_{0} e^{-\theta_{\min } \alpha}, \tag{2.99}
\end{equation*}
$$

then for some $C_{2}>0$ independent of $T, x$ and $\varphi$ it holds

$$
\left|\gamma_{x}(\varphi, T)\right| \leqslant C_{2} \max \left\{T, T^{\frac{\theta_{\min }}{\theta_{\text {max }}}+s^{\theta_{\text {max }}}}\right\}\|\varphi\|_{W_{P}^{s, t, q}} .
$$

(ii) If for some $\tilde{\varphi} \in W_{p}^{s, t, q}$ it holds for all $\alpha \geqslant 0$, for some $\lambda \in \mathbb{R}, c \geqslant 0$ and some $C=C(\lambda, c, \tilde{\varphi})>0$

$$
\begin{equation*}
\left\|\partial_{\rho} \tau(0,-\alpha, \cdot) \cdot \tilde{\varphi} \circ g_{-\alpha}\right\|_{W_{P}^{s, t, q}} \leqslant C e^{\lambda \alpha} \max \left\{1,|\alpha|^{c}\right\}, \tag{2.100}
\end{equation*}
$$

then there exists $C_{3}=C_{3}(\lambda, c)>0$ independent of $T, x$ and $\tilde{\varphi}$ such that

$$
\left|\gamma_{w, x}(\widetilde{\varphi})\right| \leqslant C C_{3} \begin{cases}T^{\frac{\lambda}{h_{\text {top }}}}(\max \{1, \log T\})^{c}, & \text { if } \lambda>0 \\ \min \{1, T\}(\max \{1, \log T\})^{c+1}, & \text { if } \lambda=0 \\ \min \{1, T\}(\max \{1, \log T\})^{c}, & \text { if } \lambda<0\end{cases}
$$

Moreover, if the bound in (2.100) holds for all $\alpha \in \mathbb{R}$ with $\lambda>0$ then

$$
\left|\gamma_{x}(\widetilde{\varphi}, T)-\gamma_{w, x}(\widetilde{\varphi})\right| \leqslant C C_{3} .
$$

Proof. Let $x \in M, T>0,0<\epsilon \leqslant \frac{1}{4}$. We define $\beta_{k}^{+}, \beta_{k}^{-} \in \mathbb{R}$ for every $k \in \mathbb{N}$ by

$$
\begin{array}{rll}
\tau\left(T, \beta_{0}^{+}, x\right)=\frac{1}{\epsilon} & \text { and } & \tau\left(\tau\left(\frac{1}{\epsilon},-\beta_{k}^{+}, g_{\beta_{k}^{+}}(x)\right), \beta_{k-1}^{+}, x\right)=1 \\
\beta_{0}^{-}:=\beta_{0}^{+} & \text {and } & \tau\left(\tau\left(-\frac{1}{\epsilon},-\beta_{k}^{-}, g_{\beta_{k}^{-}} \circ h_{T}(x)\right), \beta_{k-1}^{-}, h_{T}(x)\right)=-1 . \tag{2.101}
\end{array}
$$

If $T>1$ we assume $\beta_{0}^{+}>0$, if $T=1$ we assume $\beta_{0}^{+}=0$ and if $T<1$ we assume $\beta_{0}^{+}<0$. This is justified since $\tau(T, 0, x)=T$ and by Proposition 2.5.13 (xi). Since $\epsilon<1$ we may assume without loss of generality for all $k \in \mathbb{N}$

$$
\beta_{k}^{+}<\beta_{k-1}^{+} \quad \text { and } \quad \beta_{k}^{-}<\beta_{k-1}^{-} .
$$

Combining the definitions in (2.101) with (iii) and (xi) in Proposition 2.5.13, we find $C_{1} \geqslant 1$ independent of $\epsilon, x, k$ and $T$ such that for all $k \in \mathbb{N}$ it holds

$$
\begin{equation*}
C_{1}^{-1} \epsilon \leqslant e^{h_{\text {top }}\left(\beta_{k}^{+}-\beta_{k-1}^{+}\right)}, e^{h_{\text {top }}\left(\beta_{k}^{-}-\beta_{k-1}^{-}\right)} \leqslant C_{1} \epsilon . \tag{2.102}
\end{equation*}
$$

If $\beta_{0}^{+} \geqslant 0$ it follows for all $k \in \mathbb{Z}_{\geqslant 0}$, using the upper bounds in (2.102) and Proposition 2.5.13 (xii) on $\tau\left(T, \beta_{0}^{+}, x\right)=\frac{1}{\epsilon}$,

$$
\begin{equation*}
T\left(C_{1}^{-1} \epsilon\right)^{k+1} \leqslant e^{h_{\text {top }} \beta_{k}^{+}} \leqslant T\left(C_{1} \epsilon\right)^{k+1} . \tag{2.103}
\end{equation*}
$$

If $\beta_{0}^{+}<0$ it holds for all $k \in \mathbb{Z}_{\geqslant 0}$

$$
\begin{equation*}
\left(C_{0}^{-1} T\right)^{\frac{h_{\text {top }}}{\theta_{\mathrm{min}}}}\left(C_{1}^{-1} \epsilon\right)^{k} \leqslant e^{h_{\mathrm{top}} \beta_{k}^{+}} \leqslant\left(C_{0} T\right)^{\frac{h_{\text {top }}}{\theta_{\max }}}\left(C_{1} \epsilon\right)^{k}, \tag{2.104}
\end{equation*}
$$

where $C_{0}, \theta_{\min }$ and $\theta_{\max }$ are from the assumptions in (2.99). By symmetry we obtain analogous bounds for $\beta_{k}^{-}$. We let $w^{+}, w^{-} \in C^{\infty}(\mathbb{R},[0,1])$ such that

$$
w^{-}=w^{+} \circ(T-\cdot), \quad w_{\left\lvert\,\left(\frac{1}{2 \epsilon}, \infty\right)\right.}^{+} \equiv 1 \quad \text { and } \quad w_{\left\lvert\,\left(-\infty, \frac{1}{4 \epsilon}\right)\right.}^{+} \equiv 0 .
$$

We set

$$
w_{0}:=w^{+} \circ \tau\left(\cdot, \beta_{0}^{+}, x\right) \cdot w^{-} \circ\left(T+\tau\left(\cdot-T, \beta_{0}^{-}, h_{T}(x)\right)\right),
$$

and we set for all $k \in \mathbb{N}$

$$
\begin{aligned}
w_{k}^{+} & :=w^{+} \circ \tau\left(\cdot, \beta_{k}^{+}, x\right)-w^{+} \circ \tau\left(\cdot, \beta_{k-1}^{+}, x\right), \\
w_{k}^{-} & :=w^{-} \circ\left(T+\tau\left(\cdot-T, \beta_{k}^{-}, h_{T}(x)\right)\right)-w^{-} \circ\left(T+\tau\left(\cdot-T, \beta_{k-1}^{-}, h_{T}(x)\right)\right), \\
w_{k} & :=w_{k}^{+}+w_{k}^{-} .
\end{aligned}
$$

We let $N \in \mathbb{Z}$ for now be arbitrary. If $N \geqslant 0$ we set

$$
\begin{aligned}
w:= & \sum_{k=0}^{N} w_{k}=w_{0}+w^{+} \circ \tau\left(\cdot, \beta_{N}^{+}, x\right)-w^{+} \circ \tau\left(\cdot, \beta_{0}^{+}, x\right) \\
& +w^{-} \circ\left(T+\tau\left(\cdot-T, \beta_{N}^{-}, h_{T}(x)\right)\right)-w^{-} \circ\left(T+\tau\left(\cdot-T, \beta_{0}^{-}, h_{T}(x)\right)\right) .
\end{aligned}
$$

If $N<0$ we put $w \equiv 0$. Since $\epsilon \leqslant \frac{1}{4}$, it follows directly from the definitions of

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$\beta_{k}^{+}$and $\beta_{k}^{-}$in (2.101) that for all $k \in \mathbb{N}$

$$
\begin{aligned}
\tau\left(\frac{1}{2 \epsilon},-\beta_{k}^{+}, g_{\beta_{k}^{+}}(x)\right) & \leqslant \tau\left(\frac{1}{4 \epsilon},-\beta_{k-1}^{+}, g_{\beta_{k-1}^{+}}(x)\right) \quad \text { and } \\
\tau\left(\frac{1}{-2 \epsilon},-\beta_{k}^{-}, g_{\beta_{k}^{-}} \circ h_{T}(x)\right) & \geqslant \tau\left(-\frac{1}{4 \epsilon},-\beta_{k-1}^{-}, g_{\beta_{k-1}^{+}} \circ h_{T}(x)\right) .
\end{aligned}
$$

Together with the assumptions on the supports of $w^{+}$and $w^{-}$, we find if $N \geqslant 0$

$$
\begin{aligned}
\operatorname{supp}\left(1_{\mid[0, T]}-w\right) \subseteq & \left(0, \tau\left(\frac{1}{2 \epsilon},-\beta_{N}^{+}, g_{\beta_{N}^{+}}(x)\right)\right) \\
& \cup\left(T+\tau\left(-\frac{1}{2 \epsilon},-\beta_{N}^{-}, g_{\beta_{N}^{-}} \circ h_{T}(x)\right), T\right) .
\end{aligned}
$$

We put

$$
N:=\left\lfloor-\log \left(C_{1}^{-1} T\right) / \log \left(C_{1} \epsilon\right)\right\rfloor .
$$

Hence if $N<0$ then $T$ is bounded and if $N \geqslant 0$ then $\beta_{N}^{+}, \beta_{N}^{-} \geqslant 0$. The latter follows from the lower bounds in (2.103). Therefore the first statement in Claim (i) follows immediately, using in addition Proposition 2.5.13 (xii) and the upper bounds in (2.103).
Using Proposition 2.5.13 (iii) and also Proposition 2.5.13 (iv) in the last equality for $\tilde{w}_{0}$, we find for all $k \in \mathbb{N}$

$$
\begin{aligned}
\tilde{w}_{k}^{+} & :=w_{k}^{+} \circ \tau\left(\cdot,-\beta_{k}^{+}, g_{\beta_{k}^{+}}(x)\right)=w^{+}-w^{+} \circ \tau\left(\cdot, \beta_{k-1}^{+}-\beta_{k}^{+}, g_{\beta_{k}^{+}}(x)\right) \\
\widetilde{w}_{k}^{-} & :=w_{k}^{-}(\cdot+T) \circ \tau\left(\cdot,-\beta_{k}^{-}, g_{\beta_{k}^{-}} \circ h_{T}(x)\right) \\
& =w^{-}(T+\cdot)-w^{-} \circ\left(T+\tau\left(\cdot, \beta_{k-1}^{-}-\beta_{k}^{-}, g_{\beta_{k}^{-}} \circ h_{T}(x)\right)\right) \\
\widetilde{w}_{0} & :=w_{0} \circ \tau\left(\cdot,-\beta_{0}^{+}, g_{\beta_{0}^{+}}(x)\right) \\
& =w^{+} \cdot w^{-} \circ\left(T+\tau\left(\tau\left(\cdot,-\beta_{0}^{+}, g_{\beta_{0}^{+}}(x)\right)-T, \beta_{0}^{+}, h_{T}(x)\right)\right) \\
& =w^{+} \cdot w^{-} \circ\left(T-\frac{1}{\epsilon}+\cdot\right)
\end{aligned}
$$

For this construction it holds for all $k \in \mathbb{N}$

$$
\begin{equation*}
\operatorname{supp} \tilde{w}_{k}^{+},-\operatorname{supp} \tilde{w}_{k}^{-}, \operatorname{supp} \tilde{w}_{0} \subseteq\left[0, \frac{C_{1}^{2}}{2} 1 / \epsilon^{2}\right] \tag{2.105}
\end{equation*}
$$

Since $\partial_{\rho} \tau(\cdot, \alpha, x) \in C^{r-1}$ for all $\alpha \in \mathbb{R}$ and all $x \in M$ it follows for some constant $C_{3}>0$ for all $k \in \mathbb{N}$, all $x \in M$ and all $T>0$, using Proposition 2.5.13 (viii) and
the bounds in (2.102),

$$
\begin{equation*}
\left\|\tilde{w}_{k}^{+}\right\|_{C^{r}},\left\|\tilde{w}_{k}^{-}\right\|_{C^{r}},\left\|\tilde{w}_{0}\right\|_{C^{r}} \leqslant C_{3} . \tag{2.106}
\end{equation*}
$$

We note

$$
\gamma_{w, x}(\varphi)=\gamma_{w o(\cdot+T), h_{T}(x)}(\varphi) .
$$

Assuming $N \geqslant 0$, together with the equality in (2.97), we find the local decomposition $^{6}$ for all $\varphi \in W_{p}^{s, t, q}$

$$
\begin{align*}
& \gamma_{w, x}(\varphi)= \gamma_{\widetilde{w}_{0}, g_{\beta_{0}}}(x) \\
&\left.+\mathcal{L}_{\beta_{0}^{+}, \phi_{\beta_{0}^{+}}} \varphi\right)  \tag{2.107}\\
&+\sum_{k=1}^{N} \gamma_{\widetilde{w}_{k}^{+}, g_{\beta_{k}^{+}(x)}}\left(\mathcal{L}_{\beta_{k}^{+}, \phi_{\beta_{k}^{+}}} \varphi\right)+\gamma_{\widetilde{w}_{k}^{-}, g_{\beta_{k}^{-}}} o_{T}(x) \\
&\left(\mathcal{L}_{\beta_{k}^{-}, \phi_{\beta_{k}^{-}}} \varphi\right) .
\end{align*}
$$

Using the bound in Lemma 2.5.12 with the bounds in (2.106), and using the assumption in (2.100) for some $\tilde{\varphi} \in W_{p}^{s, t, q}$

$$
\begin{equation*}
\left\|\mathcal{L}_{\alpha, \phi_{\alpha}} \tilde{\varphi}\right\|_{W_{p}^{s, t, q}}=\left\|\partial_{\rho} \tau(0,-\alpha, \cdot) \cdot \tilde{\varphi} \circ g_{-\alpha}\right\|_{W_{p}^{s, t, q}} \leqslant C e^{\lambda \alpha} \max \left\{1,|\alpha|^{c}\right\}, \tag{2.108}
\end{equation*}
$$

and the bounds in (2.103), we estimate the right-hand side in the decomposition in (2.107) for some constant $C_{4}=C_{4}(c)>0$, recalling that $T$ is uniformly bounded from below if $N \geqslant 0$, and conclude the first statement in Claim (ii):

$$
\begin{equation*}
\left|\gamma_{w, x}(\varphi)\right| \leqslant C C_{4} T^{\frac{\lambda}{h_{\text {top }}}} \sum_{k=0}^{N}\left(C_{1} \epsilon\right)^{k \frac{\lambda}{h_{\text {top }}}}\left(\max \left\{1,(k+1)\left|\log \left(C_{1} \epsilon\right)\right|, \log T\right\}\right)^{c} . \tag{2.109}
\end{equation*}
$$

If $N<0$ then $w \equiv 0$ and $T$ is uniformly bounded from above and we conclude as well. To see the second statements in Claims (i)-(ii), we recall that the construction of the functions $w_{k}$ is valid for every $T>0$ and hence

$$
\sum_{k=0}^{\infty} w_{k}=1_{\mid(0, T)} .
$$

Since for all $\varphi \in C_{X}^{r-1}(M)$ it holds

$$
\gamma_{x}(\varphi, T)-\gamma_{\sum_{k=0}^{\infty} w_{k}, x}(\varphi)=0
$$

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we find by density for all $\varphi \in W_{p}^{s, t, q}$

$$
\gamma_{x}(\varphi, T)=\gamma_{\sum_{k=0}^{\infty} w_{k}, x}(\varphi) .
$$

It holds

$$
\operatorname{supp} w_{0} \subseteq[0, T] .
$$

Comparing with the supports in (2.105), together with the bounds in (2.103) and (2.104), we find some $C_{5}>0$ independent of $k, T, x$ and $\varphi$ such that for all $k \in \mathbb{N}$, if $\beta_{k}^{+} \geqslant 0$ and $\beta_{k}^{-} \geqslant 0$,

$$
\left(T-\operatorname{supp} w_{k}^{-}\right), \operatorname{supp} w_{k}^{+} \subseteq\left[0, C_{5}\left(C_{1} \epsilon\right)^{k} T\right],
$$

respectively, and if $\beta_{k}^{+}<0$ and $\beta_{k}^{-}<0$,

$$
\left(T-\operatorname{supp} w_{k}^{-}\right), \operatorname{supp} w_{k}^{+} \subseteq\left[0, C_{5} T^{\frac{\theta_{\min }}{\theta_{\max }}}\left(\left(C_{1} \epsilon\right)^{k}\right)^{\frac{\theta_{\min }}{h_{\mathrm{top}}}}\right] .
$$

Moreover, we find for some constant $C_{5} \geqslant 1$ for all $\rho_{1}, \rho_{2} \in \mathbb{R}$, all $\alpha,-s \geqslant 0$ and all $x \in M$, using Proposition 2.5.13 (iv) and the assumption of the upper bound for $\partial_{\rho} \tau$ in (2.99),
$\left(\tau\left(\rho_{1}, \alpha, x\right)-\tau\left(\rho_{2}, \alpha, x\right)\right)^{-s}=\tau\left(\rho_{1}-\rho_{2}, \alpha, h_{\rho_{2}}(x)\right)^{-s} \leqslant C_{5}^{-s}\left(\rho_{1}-\rho_{2}\right)^{-s} e^{s \theta_{\min } \alpha}$.

If $\alpha \leqslant 0$ it holds analogously, now using the lower bound for $\partial_{\rho} \tau$ in (2.99),
$\left(\tau\left(\rho_{1}, \alpha, x\right)-\tau\left(\rho_{2}, \alpha, x\right)\right)^{-s}=\tau\left(\rho_{1}-\rho_{2}, \alpha, h_{\rho_{2}}(x)\right)^{-s} \leqslant C_{5}^{-s}\left(\rho_{1}-\rho_{2}\right)^{-s} e^{s \theta_{\max } \alpha}$.

Since $0<-s<1$ and $0<\theta_{\min } \leqslant \theta_{\max }$ it holds for some constant $C_{6}>0$ independent of $T, x$ and $\varphi$, using the lower bounds in (2.103)-(2.104), for all $k \in \mathbb{N}$ and for all $T>0$,

$$
\left\|w_{k}\right\|_{C^{-s}} \leqslant C_{6} \max \left\{1, T^{s \frac{\theta_{\max }}{\theta_{\min }}}\left(\left(C_{1} \epsilon\right)^{k}\right)^{s^{\theta_{\max }}}\right\} .
$$

Then we estimate for every $\varphi \in W_{p}^{s, t, q}$, using Lemma 2.5.12 and $-s<\frac{\theta_{\text {min }}}{\theta_{\max }}$, for
some constants $C_{7}, C_{8}>0$ independent of $T, x, w$ and $\varphi$

$$
\begin{aligned}
& \left|\gamma_{x}(\varphi, T)\right| \leqslant \lim _{n \rightarrow \infty}\left|\sum_{k=0}^{n} \gamma_{w_{k}, x}(\varphi)\right| \\
& \leqslant C_{7} \sum_{k=0}^{\infty} \max \left\{T, T^{\frac{\theta_{\min }}{\theta_{\max }}+s \frac{\theta_{\max }}{\theta_{\min }}}\right\}\left(\left(C_{1} \epsilon\right)^{k}\right)^{\frac{\theta_{\min }}{h_{\text {top }}}} \max \left\{1,\left(\left(C_{1} \epsilon\right)^{k}\right)^{s \frac{\theta_{\max }}{\mathrm{m}_{\text {top }}}}\right\}\|\varphi\|_{W_{p}^{s, t, q}} \\
& \leqslant C_{8} \max \left\{T, T^{\frac{\theta_{\min }}{\theta_{\max }}+s s^{\theta_{\max }}} \theta_{\min }\right.
\end{aligned}\|\varphi\|_{W_{p}^{s, t, q}} . \quad .
$$

This yields the second statement in Claim (i). On the other hand, using the equality in (2.98) and assuming $N>0$, we find,

$$
\begin{aligned}
& \gamma_{x}(\tilde{\varphi}, T)-\sum_{k=0}^{N} \gamma_{w_{k}, x}(\tilde{\varphi})=\sum_{k=N+1}^{\infty} \gamma_{\widetilde{w}_{k}^{+}, g_{\beta_{k}^{+}(x)}}\left(\partial_{\rho} \tau\left(0,-\beta_{k}^{+}, \cdot\right) \cdot \tilde{\varphi} \circ g_{-\beta_{k}^{+}}\right) \\
& +\sum_{N+1}^{\infty} \gamma_{\widetilde{w}_{k}^{-}, g_{\beta_{k}^{-}} \circ h_{T}(x)}\left(\partial_{\rho} \tau\left(0,-\beta_{k}^{-}, \cdot\right) \cdot \tilde{\varphi} \circ g_{-\beta_{k}^{-}}\right) .
\end{aligned}
$$

Then we proceed analogously as for the bound in (2.109), now using the upper bounds in (2.103) and the assumption in (2.108) for all $\alpha \in \mathbb{R}$ and some $\lambda>0$, $c \geqslant 0$ (recall that $T\left(C_{1} \epsilon\right)^{N}$ is bounded from above). If $N \leqslant 1$ then $T^{\lambda}|\log T|^{c}$ is bounded from above and we conclude as well, now using the upper bounds in (2.104).

Remark 2.5.15. The second statement in Lemma 2.5.14 (i) can be used to avoid the $\|\varphi\|_{C^{0}}$-term in the bound of the error term in Theorem 2.5.7. However the required range for s may not be very large (except in the case of constant vector fields). The second statement in Lemma 2.5.14 (ii) is free from an additional condition on $s$. We use it in the following subsection in the proof of Lemma 2.5.18 (v) and Theorem 2.5.7. Both statements give also bounds for all values $T>0$ which seems to be new.

### 2.5.3 Showing $\lambda_{\text {max }}=h_{\text {top }}$ and Theorem 2.5.7

In this subsection we shall prove Theorem 2.5.7. First, we state and prove two lemmas which will imply that $\lambda_{\max }=h_{\mathrm{top}}$, assuming $\lambda_{\min }<\lambda_{\max }$, is a simple eigenvalue and that $\lambda_{\text {max }}$ is uniquely attained.
We remind the reader that uniqueness and simplicity of the spectral bound is known to hold (see [22, Lemma 5.1], [21]) for the spectrum of mixing Anosov flows (which are not necessarily contact), but for different anisotropic spaces, and only for the potential $V$ given by the Jacobian of the flow (and associated
to the SRB measure).
For the sake of the next two lemmas we have to introduce the following condition ${ }^{7}$ :

Condition 2.5.16 (Strong vanishing). Let $0<t, q,-s<r-1$ and let $p \in[1, \infty]$.
Let $\varphi_{\alpha} \in W_{p}^{s, t, q}$ for all $\alpha \geqslant 0$ such that $\left\|\varphi_{\alpha}\right\|_{W_{p}^{s, t, q}}=1$ and

$$
\limsup _{\alpha \rightarrow \infty} e^{-h_{\text {top }} \alpha}\left\|\mathcal{L}_{\alpha, \phi_{\alpha}} \varphi_{\alpha}\right\|_{W_{p}^{s, t, q}}>0 .
$$

If for some $-s \leqslant-s^{\prime}$ for all $x \in M$ and all $w \in C_{0}^{-s^{\prime}}(\mathbb{R})$

$$
\lim _{\alpha \rightarrow \infty} \gamma_{w, x}\left(\frac{\mathcal{L}_{\alpha, \phi_{\alpha}} \varphi_{\alpha}}{\left\|\mathcal{L}_{\alpha, \phi_{\alpha}} \varphi_{\alpha}\right\|_{W_{p}^{s, t, q}}}\right)=0
$$

then $\lim _{\alpha \rightarrow \infty}\left\|\varphi_{\alpha}\right\|_{W_{p}^{s, t, q}}=0$.
We give the upper bound on the spectral radius:
Lemma 2.5.17 (Upper bound on the spectral radius). Let $0<t, q,-s<r-1$ and let $p \in[1, \infty]$. For all $x \in M$ and all $\alpha \geqslant 0$ let $\partial_{\rho} \tau(\cdot,-\alpha, x) \in C^{r-1}(\mathbb{R}, M)$. Under Condition 2.5.16, With the choice $\phi_{\alpha}=\partial_{\rho} \tau(0,-\alpha, \cdot)$ for some constant $C>0$ it holds for all $\alpha \geqslant 0$

$$
\left\|\mathcal{L}_{\alpha, \phi_{\alpha}}\right\|_{W_{p}^{s, t, q} \rightarrow W_{p}^{s, t, q}} \leqslant C e^{h_{\text {top } \alpha} \alpha} .
$$

Proof. We show the claim on $\left\|\mathcal{L}_{\alpha, \phi_{\alpha}}\right\|_{W_{p}^{s, t, q} \rightarrow W_{p}^{s, t, q}}$ by contradiction. Suppose

$$
e^{-h_{\text {top } \alpha}}\left\|\mathcal{L}_{\alpha, \phi_{\alpha}}\right\|_{W_{p}^{s, t, q} \rightarrow W_{p}^{s, t, q}} \rightarrow \infty \quad \text { as } \quad \alpha \rightarrow \infty .
$$

Then there exists $\varphi_{\alpha} \in W_{p}^{s, t, q}$ such that $\left\|\varphi_{\alpha}\right\|_{W_{p}^{s, t, q}}=1$ and

$$
\begin{equation*}
\left\|e^{-h_{\text {top }} \alpha} \mathcal{L}_{\alpha, \phi_{\alpha}} \varphi_{\alpha}\right\|_{W_{p}^{s, t, q}} \rightarrow \infty \quad \text { as } \quad \alpha \rightarrow \infty . \tag{2.110}
\end{equation*}
$$

We assume for some $w \in C_{0}^{s}(\mathbb{R})$ and some $x \in M$

$$
\begin{equation*}
\liminf _{\alpha \rightarrow \infty}\left|\gamma_{w, x}\left(\frac{\mathcal{L}_{\alpha, \phi_{\alpha}} \varphi_{\alpha}}{\left\|\mathcal{L}_{\alpha, \phi_{\alpha}} \varphi_{\alpha}\right\|_{W_{p}^{s, t, q}}}\right)\right|>0 \tag{2.111}
\end{equation*}
$$

[^7]This assumption is justified, assuming Condition 2.5.16. We choose $T \geqslant 1$ and $\alpha \geqslant 0$ such that

$$
\tau(T, \alpha, x)=1 .
$$

Then, using Proposition 2.5.13 (xii), we find for some constant $C \geqslant 1$

$$
\begin{equation*}
C^{-1} e^{h_{\text {top }} \alpha} \leqslant T \leqslant C e^{h_{\text {top }} \alpha} . \tag{2.112}
\end{equation*}
$$

We have, using the equality given in (2.97),

$$
\begin{equation*}
\gamma_{w \circ \tau(\cdot, \alpha, x), x}\left(\varphi_{\alpha}\right)=\gamma_{w, g_{\alpha}(x)}\left(\mathcal{L}_{\alpha, \phi_{\alpha}} \varphi_{\alpha}\right) \tag{2.113}
\end{equation*}
$$

We recall $-s<r-1$. Therefore the norm $\|w \circ \tau(\cdot, \alpha, x)\|_{C^{-s}}$ is bounded as $\alpha \rightarrow \infty$, using Proposition 2.5.13 (x). By Lemma 2.5.12, the linear functionals $\gamma_{w \circ \tau(\cdot, \alpha, x), x}$ and $\gamma_{w, g_{\alpha}(x)}$ which appear in (2.113) are continuous on $W_{p}^{s, t, q}$. Hence the left-hand side in (2.113) grows at most by $T$ as $\alpha \rightarrow \infty$ uniformly in $x$.
Then, comparing with the estimates for $T$ in (2.112), using the assumption in (2.111), this contradicts the assumption in (2.110) and we conclude.

We next show the lower bound (and uniqueness and simplicity of the spectral bound $\lambda_{\text {max }}$ ):

Lemma 2.5.18 (Invariant measure and spectral bound). Let $\mu$ be the unique Borel probability measure which is invariant by the horocycle flow $h_{\rho}$. Let $p \in$ $[1, \infty]$ and let $s<0<q \leqslant t$ such that $t-s<r-1$. It holds:
(i) $\mu \in\left(W_{p}^{s, t, q}\right)^{\prime}$,
(ii) $\mathcal{L}_{\alpha, \phi_{\alpha}}^{\prime} \mu=e^{h_{\text {top }} \alpha} \mu\left(\mathcal{L}_{\alpha, \phi_{\alpha}}^{\prime}\right.$ denotes the adjoint operator of $\left.\mathcal{L}_{\alpha, \phi_{\alpha}}\right)$,
(iii) $\left.h_{\text {top }} \in \sigma(X+V)\right|_{W_{p}^{s, t, q}}$. Moreover, assuming Condition 2.5.16, it holds:
(iv) $\lambda_{\text {max }}=h_{\text {top }}$.
(v) The spectral bound $\lambda_{\text {max }}$ is uniquely attained by the simple eigenvalue $h_{\text {top }}$, assuming $\lambda_{\text {min }}<\lambda_{\text {max }}$.

The vector $\mu$ is also invariant by the adjoint horocycle flow since the time average converges to the (unique) ergodic mean (a result by Marcus [50]). This is in analogy to [31, Lemma 2.11].

Proof. We note for every $\varphi \in C_{X}^{r-1}(M)$, using [49, Theorem 2.1] for the first, the equality in (2.97) for the second and [49, Lemma 3.1] for the third equality, for some $\lambda>0$, for every $\alpha \geqslant 0$

$$
\begin{align*}
\mu(\varphi) & =\lim _{T \rightarrow \infty} \frac{1}{T} \gamma_{x}(\varphi, T)=\lim _{T \rightarrow \infty} \frac{\tau(T, \alpha, x)}{T} \frac{1}{\tau(T, \alpha, x)} \gamma_{g_{\alpha(x)}}\left(\mathcal{L}_{\alpha, \phi_{\alpha}} \varphi, \tau(T, \alpha, x)\right) \\
& =\lambda^{-\alpha} \mu\left(\mathcal{L}_{\alpha, \phi_{\alpha}} \varphi\right) . \tag{2.114}
\end{align*}
$$

To see $\lambda=e^{h_{\text {top }}}$ we refer to [49, p.84] (alternatively use Proposition 2.5.13 (xi)). Using Claims (i)-(ii) with $\lambda=h_{\text {top }}$ in Lemma 2.5.14 together with the bound given by Lemma 2.5.17, there is $w \in C^{r}(\mathbb{R})$ and a constant $C_{1}>0$ such that for all $\varphi \in C_{X}^{r-1}(M)$

$$
\begin{equation*}
|\mu(\varphi)| \leqslant \lim _{T \rightarrow \infty}\left|\frac{1}{T} \gamma_{w, x}(\varphi)\right| \leqslant C_{1}\|\varphi\|_{W_{p}^{s, t, q}} . \tag{2.115}
\end{equation*}
$$

Claim (iii) follows from $\left.\sigma\left((X+V)^{\prime}\right)\right|_{\left(W_{p}^{s, t, q}\right)^{\prime}}=\left.\sigma(X+V)\right|_{W_{p}^{s, t, q},}$ using [46, Section II.2.5]. Claim (iv) follows from (iii) together with Lemma 2.5.17. To see Claim (v), first we note that all $\lambda \in \sigma(X+V)$ such that $\Re \lambda=h_{\text {top }}$ are eigenvalues, using Lemma 2.4.10 together with the assumption $\lambda_{\min }<\lambda_{\max }$. Using Claim (iii), there exists $\mathcal{D}_{1} \in W_{p}^{s, t, q}$ such that $\mathcal{L}_{\alpha, \phi_{\alpha}} \mathcal{D}_{1}=e^{h_{\text {top }} \alpha} \mathcal{D}_{1}$ for all $\alpha \geqslant 0$. We let $\mathcal{D}_{1} \neq \mathcal{D}_{2} \in W_{p}^{s, t, q} \backslash\{0\}$ such that $\mathcal{L}_{\alpha, \phi_{\alpha}} \mathcal{D}_{2}=e^{\lambda \alpha} \mathcal{D}_{2}$ for all $\alpha \geqslant 0$, where $\lambda \in \mathbb{C}$ and $\Re \lambda=h_{\text {top }}$. Then it holds, using Claim (ii) for the last equality,

$$
e^{\lambda \alpha} \mu\left(\mathcal{D}_{2}\right)=\mu\left(\mathcal{L}_{\alpha, \phi_{\alpha}} \mathcal{D}_{2}\right)=e^{h_{\text {top } \alpha}} \mu\left(\mathcal{D}_{2}\right) .
$$

Since $\lambda \neq h_{\text {top }}$ it holds $\mu\left(\mathcal{D}_{2}\right)=0$. In fact, by same reasoning we can always assume $\mu\left(\mathcal{D}_{2}\right)=0$ if $\lambda \neq h_{\text {top }}$. And if $\lambda=h_{\text {top }}$ there are only finitely many such $\mathcal{D}_{2}$ and we can again assume $\mu\left(\mathcal{D}_{2}\right)=0$ by a change of basis. The upshot is that the following reasoning works always if $\Re \lambda \geqslant h_{\text {top }}$ and $\mu\left(\mathcal{D}_{2}\right)=0$.
Then, using Claim (i) and the equality in (2.114), for every $\epsilon>0$ there is $\varphi \in C_{X}^{r-1}(M)$ such that for all $\alpha \in \mathbb{R}$ and for all $x \in M$

$$
\begin{equation*}
\lim _{T \rightarrow \infty}\left|\frac{1}{T} \gamma_{g-\alpha}(x)(\varphi, T)\right|=|\mu(\varphi)| \leqslant \epsilon . \tag{2.116}
\end{equation*}
$$

Using Lemma 2.5.14 (i), for all $T>0$, for all $x \in M$ and for all $\alpha \in \mathbb{R}$ there exists $w \in C^{r}$ such that

$$
\begin{equation*}
\lim _{T \rightarrow \infty}\left|\frac{1}{T} \gamma_{g_{-\alpha}(x)}(\varphi, T)\right|=\lim _{T \rightarrow \infty}\left|\frac{1}{T} \gamma_{w, g_{-\alpha}(x)}(\varphi)\right| . \tag{2.117}
\end{equation*}
$$

Since $\partial_{\rho} \tau(0,-\alpha, \cdot) \in C^{r-1}$ for all $\alpha \in \mathbb{R}$, using Lemma 2.3.3, we find $\mathcal{L}_{\alpha, \phi_{\alpha}} \mathcal{D}_{2}=$ $e^{\lambda \alpha} \mathcal{D}_{2}$ for all $\alpha \in \mathbb{R}$ which matches the condition (2.100) in Lemma 2.5.14 (ii). Then, using Lemma 2.5.14 (ii) for the upper bound and the equality in (2.97) for the last step, we find for some constant $C_{2}>0$ independent of $x, \alpha$ and $\varphi$, for all $\alpha \in \mathbb{R}$

$$
\begin{align*}
\lim _{T \rightarrow \infty}\left|\frac{1}{T} \gamma_{w, g_{-\alpha}(x)}(\varphi)\right| & =\lim _{T \rightarrow \infty}\left|\frac{1}{T} \gamma_{w, g_{-\alpha}(x)}\left(\mathcal{D}_{2}\right)+\frac{1}{T} \gamma_{w, g_{-\alpha}(x)}\left(\varphi-\mathcal{D}_{2}\right)\right| \\
& \geqslant \lim _{T \rightarrow \infty}\left|\frac{1}{T} \gamma_{g_{-\alpha}(x)}\left(\mathcal{D}_{2}, T\right)\right|-C_{2}\left\|\varphi-\mathcal{D}_{2}\right\|_{W_{p}^{s, t, q}} \\
& =\lim _{T \rightarrow \infty}\left|\frac{e^{h_{\text {top }} \alpha}}{T} \gamma_{x}\left(\mathcal{D}_{2}, \tau(T, \alpha, x)\right)\right|-C_{2}\left\|\varphi-\mathcal{D}_{2}\right\|_{W_{p}^{s, t, q}} \tag{2.118}
\end{align*}
$$

By density of $C_{X}^{r-1}(M)$ in $W_{p}^{s, t, q}$ we assume

$$
\left\|\varphi-\mathcal{D}_{2}\right\|_{W_{P}^{s, t, q}} \leqslant \epsilon
$$

For every $T \geqslant T_{0} \geqslant 1$ we let $\alpha \geqslant 0$ such that $\tau(T, \alpha, x)=T_{0}$. By Proposition 2.5.13 there is $C_{3} \geqslant 1$ independent of $T$ and $x$ such that $e^{h_{\mathrm{top}} \alpha} \geqslant C_{3}^{-1} \frac{T}{T_{0}}$. Since $\epsilon>0$ was arbitrary we conclude for all $T_{0} \geqslant 1$ and all $x \in M$, using the estimates (2.116)-(2.118),

$$
\gamma_{x}\left(\mathcal{D}_{2}, T_{0}\right)=0
$$

On the other hand we find for all $T_{1}, T_{0} \geqslant 1$

$$
\gamma_{x}\left(\mathcal{D}_{2}, T_{0}-T_{1}\right)=\gamma_{h_{-T_{1}}(x)}\left(\mathcal{D}_{2}, T_{0}\right)-\gamma_{h_{-T_{1}}(x)}\left(\mathcal{D}_{2}, T_{1}\right)=0
$$

Hence it holds $\gamma_{x}\left(\mathcal{D}_{2}, T\right)=0$ for every $T \in \mathbb{R}$ and every $x \in M$. Then for every $w \in C_{0}^{s+1}$ we find, using integration by parts,

$$
\gamma_{w, x}\left(\mathcal{D}_{2}\right)=-\int_{\mathbb{R}}\left(\partial_{\rho} w\right)(\rho) \cdot \gamma_{x}\left(\mathcal{D}_{2}, \rho\right) \mathrm{d} \rho
$$

Since $\gamma_{x}\left(\mathcal{D}_{2}, \rho\right)=0$ for all $\rho \geqslant 0$ we conclude $\gamma_{w, x}\left(\mathcal{D}_{2}\right)=0$. Then, using Condition 2.5.16, we find $\mathcal{D}_{2} \equiv 0$ but we assumed $\mathcal{D}_{2} \not \equiv 0$.

Proof of Theorem 2.5.7. By assumption

$$
\max \left\{\lambda_{\min }, 0\right\}<\delta \leqslant \lambda_{\max }=h_{\mathrm{top}} .
$$

We note that we have always $\lambda_{\max }=h_{\text {top }}$ and uniqueness and simplicity of $\lambda_{\max }$
under Condition 2.5.16. Using the equality in (2.67) for the projectors $\Pi_{\lambda, i}$, we have for all $1 \leqslant i \leqslant n_{\lambda}$

$$
\Pi_{\lambda, i} \varphi=\sum_{j=1}^{m_{\lambda, i}} \mathcal{O}_{(\lambda, i, j)}(\varphi) \mathcal{D}_{(\lambda, i, j)}
$$

Recalling the nil-potent operators $\mathcal{N}_{\lambda, i}$ of finite rank (e.g. see in (2.66)), using the formula for the matrix action $\mathcal{L}_{\alpha, \phi_{\alpha}} \Pi_{\lambda, i}=\exp (\lambda \alpha) \exp \left(\mathcal{N}_{\lambda, i} \alpha\right) \Pi_{\lambda, i}$ for all $\alpha \geqslant 0$ and

$$
\exp (-\lambda \alpha) \exp \left(-\mathcal{N}_{\lambda, i} \alpha\right) \mathcal{L}_{\alpha, \phi_{\alpha}} \Pi_{\lambda, i}=\Pi_{\lambda, i}
$$

we find for some constant $C_{1}=C_{1}(\lambda, i, j)$ for all $\alpha \in \mathbb{R}$

$$
\left\|\partial_{\rho}(0,-\alpha, \cdot) \cdot \mathcal{D}_{(\lambda, i, j)} \circ g_{-\alpha}\right\|_{W_{p}^{s, t, q}} \leqslant C_{1} \exp (\Re \lambda \alpha) \max \left\{1,|\alpha|^{j-1}\right\} .
$$

Hence $\mathcal{D}_{(\lambda, i, j)}$ satisfies the upper bound in (2.100) for all $\alpha \in \mathbb{R}$ if $\Re \lambda>0$.
Inspecting the end of the proof of Lemma 2.5.18, we notice that all eigendistributions $\mathcal{D}_{(\lambda, i, j)}$ associated to some eigenvalue $\lambda$ with $\Re \lambda \geqslant h_{\text {top }}$ do not contribute to the expansion except $\mathcal{D}_{\left(h_{\text {top }}, 1,1\right)}$. This follows, if $j=1$ using that $\gamma_{x}\left(\mathcal{D}_{(\lambda, i, 1)}, T\right)=0$ for all $T \geqslant 0$ and all $x \in M$. If $j>1$ we arrive at the same conclusion, using in the estimate in (2.118) for all $\alpha \geqslant 0$

$$
\mathcal{L}_{\alpha, \phi_{\alpha}} \mathcal{D}_{(\lambda, i, j)}=\exp (\lambda \alpha) \exp \left(\mathcal{N}_{\lambda, i} \alpha\right) \mathcal{D}_{(\lambda, i, j)}
$$

Let $\lambda \in \Sigma_{\delta}=\left.\sigma(X+V)\right|_{W_{p}^{s, t, q}} \cap\{z \in \mathbb{C} \mid \Re z \geqslant \delta\}$. For every $T \geqslant 0$ and every $x \in M$ we set, using $w \in C^{r}$ given in Lemma 2.5.14,

$$
c_{(\lambda, i, j)}=c_{(\lambda, i, j)}(T, x):=T^{-\frac{\lambda}{h_{\text {top }}}} \max \left\{1,|\log T|^{1-j}\right\} \gamma_{w, x}\left(\mathcal{D}_{(\lambda, i, j)}\right) .
$$

Then, using the first statement in Lemma 2.5.14 (ii), the coefficients $c_{(\lambda, i, j)}$ are bounded independently of $T$ and $x$. It holds

$$
\begin{aligned}
\gamma_{w, x}\left(\Pi_{\lambda, i} \varphi\right) & =\sum_{j=1}^{m_{\lambda, i}} \mathcal{O}_{(\lambda, i, j)}(\varphi) \gamma_{w, x}\left(\mathcal{D}_{(\lambda, i, j)}\right) \\
& =\sum_{j=1}^{m_{\lambda, i}} c_{(\lambda, i, j)} T^{\frac{\lambda}{n_{\text {top }}}} \max \left\{1,|\log T|^{j-1}\right\} \mathcal{O}_{(\lambda, i, j)}(\varphi) .
\end{aligned}
$$

We let $\mu$ as given in Lemma 2.5.18. Using Lemma 2.5.18 (v), and assuming
$T \geqslant e$, we find for every finite subset $\Lambda_{\delta} \subseteq \Sigma_{\delta}$

$$
\gamma_{x}(\cdot, T)=\gamma_{x}\left(\mathcal{D}_{\left(h_{\text {top }}, 1,1\right)}, T\right) \mu+\sum_{\substack{\lambda \in \Lambda_{\delta} \\ \Re \lambda<h_{\text {top }}}} \sum_{i=1}^{n_{\lambda}} \sum_{j=1}^{m_{\lambda, i}} c_{(\lambda, i, j)} T^{\frac{\lambda}{h_{\text {top }}}}(\log T)^{j-1} \mathcal{O}_{(\lambda, i, j)}+\mathcal{E}_{T, x, \Lambda_{\delta}},
$$

where the remainder term is

$$
\begin{align*}
\mathcal{E}_{T, x, \Lambda_{\delta}}:= & \left(\gamma_{w, x}\left(\mathcal{D}_{h_{\text {top }}, 1,1}\right)-\gamma_{x}\left(\mathcal{D}_{\left(h_{\text {top }, 1,1)}, T\right)}\right) \mu\right. \\
& +\gamma_{w, x}\left(\mathrm{id}-\sum_{\lambda \in \Lambda_{\delta}} \sum_{i=1}^{n_{\lambda}} \Pi_{\lambda, i}\right)+\left(\gamma_{x}(\cdot, T)-\gamma_{w, x}\right) . \tag{2.119}
\end{align*}
$$

The existence of the limit $\lim _{T \rightarrow \infty} T^{-1} \gamma_{x}\left(\mathcal{D}_{\left(h_{\text {top }}, 1,1\right)}, T\right)$ is shown by analogue estimates (2.116)-(2.118). Then the statement on the limit $\lim _{T \rightarrow \infty} T^{-1} \mathcal{E}_{T, x, \Lambda_{\delta}}(\varphi)$ follows, using unique ergodicity of the horocycle flow [49, Theorem 2.1] and finiteness of $\Lambda_{\delta}$. We bound $\left|\mathcal{E}_{T, x, \Lambda_{\delta}}(\varphi)\right|$ as required, using the first statement in Lemma 2.5.14 (i) and the full statement in Lemma 2.5.14 (ii) together with the assumed upper bound in (2.88).
The additional claims under Condition 2.4.11 can be seen as follows (see also the remarks above and below Condition 2.4.11): The finiteness of $\Sigma_{\delta}$ follows from [19, Theorem 1]. To this end we have to show that [19, Assumption 1-3A] are satisfied for the renormalized semigroup $e^{-h_{\text {top }} \alpha} \mathcal{L}_{\alpha, \phi_{\alpha}}: W_{p}^{s, t, q} \rightarrow W_{p}^{s, t, q}$. In fact Condition 2.4.11 yields just a reformulation of [19, Assumption 3A] for the resolvent of the generator $X+V-h_{\text {top }}$. Now [19, Assumption 1] states that for some Banach space $W_{p}^{s, t, q} \subset \mathcal{B}$ it holds

$$
\begin{equation*}
\sup _{\alpha \geqslant 0} \frac{1}{\alpha}\left\|\operatorname{id}-e^{-h_{\text {top } \alpha}} \mathcal{L}_{\alpha, \phi_{\alpha}}\right\|_{W_{p}^{s, t, q} \rightarrow \mathcal{B}}<\infty . \tag{2.120}
\end{equation*}
$$

We set $\mathcal{B}:=W_{p}^{s, t, q-1}$. We bound the left-hand side in (2.120), using the equality in (2.69) together with Lemma 2.3.3, Lemma 2.3.5 and Lemma 2.5.17. Now [19, Assumption 2] just states that the essential spectral bound of $X+V-h_{\text {top }}$ is bounded by some $\lambda<0$, where $V=-\partial_{\alpha} \partial_{\rho} \tau(0,0, \cdot)$. By assumption it holds $\lambda \leqslant \lambda_{\min }-h_{\text {top }}<0$. Finally, the claimed choice $c=\delta+\epsilon$ for all $\epsilon>0$ follows from [19, Theorem 1] as well. In particular, this choice for $c$ follows if for all $\alpha \geqslant 0$ and for all $\epsilon>0$ there exists $C_{2}=C_{2}(\delta, \epsilon, \varphi)$ such that

$$
\left\|\mathcal{L}_{\alpha, \phi_{\alpha}}\left(\mathrm{id}-\sum_{\lambda \in \Sigma_{\delta}} \sum_{i=1}^{n_{\lambda}} \Pi_{\lambda, i}\right) \varphi\right\|_{W_{p}^{s, t, q}} \leqslant C_{2} e^{(\delta+\epsilon) \alpha}
$$

2 Horocycle averages on closed manifolds

We set $\varphi_{\delta}:=\varphi-\sum_{\lambda \in \Sigma_{\delta}} \sum_{i=1}^{n_{\lambda}} \Pi_{\lambda, i} \varphi$. If $t-q+1<r-1$ it follows, using Lemma 2.3.3, Lemma 2.3.5 and [19, Theorem 1], for some constants $C_{3}, C_{4}=C_{4}(\epsilon)>0$

$$
\begin{aligned}
&\left\|\mathcal{L}_{\alpha, \phi_{\alpha}} \varphi_{\delta}\right\|_{W_{p}^{s, t, q}} \leqslant C_{3}\left\|\mathcal{L}_{\alpha, \phi_{\alpha}}(X+V) \varphi_{\delta}\right\|_{W_{p}^{s, t, q-1}}+C_{3}\left\|\mathcal{L}_{\alpha, \phi_{\alpha}} \varphi_{\delta}\right\|_{W_{p}^{s, t, q-1}} \\
& \leqslant C_{3}\left\|\mathcal{L}_{\alpha, \phi_{\alpha}}\left(X+V-h_{\mathrm{top}}\right) \varphi_{\delta}\right\|_{W_{p}^{s, t, q-1}}+\left(C_{3}+h_{\mathrm{top}}\right)\left\|\mathcal{L}_{\alpha, \phi_{\alpha}} \varphi_{\delta}\right\|_{W_{p}^{s, t, q-1}} \\
& \leqslant C_{4} e^{(\delta+\epsilon) \alpha}\left(\left\|\left(X+V-h_{\mathrm{top}}\right)^{2} \varphi\right\|_{W_{p}^{s, t, q}}+\left\|\left(X+V-h_{\mathrm{top}}\right) \varphi\right\|_{W_{p}^{s, t, q}}\right) .
\end{aligned}
$$

Boundedness of the last estimate follows if $q<r-2$ because then $\|X(V \varphi)\|_{W_{p}^{s, t, q}}$ and $\left\|X^{2} \varphi\right\|_{W_{p}^{s, t, q}}$ are bounded, recalling $\varphi \in C_{X}^{r-1}(M)$ and $V \in C^{r-1}$. Combining the required bounds for $q$ yields

$$
t-r+2<q<r-2 .
$$

Since we required $q>0$ it is enough to require $t-r+2 \leqslant 0$ and $0<r-2$ which yields the additional condition on $t$ and $r$.

## Appendix A

For the readers convenience we give a proof of a well-known result:
Lemma A. 1 (Fixed points). Let $M$ be $2 \times 2$ integer matrix acting on $\mathbb{T}^{2}$. Assume that $\operatorname{det}(\operatorname{id}-M) \neq 0$. Then the following holds:
(i) The number $N_{M}$ of fixed points of $M$ is given by $N_{M}=|\operatorname{det}(\mathrm{id}-M)|$.
(ii) There exists a disjoint partition $D_{j} \subseteq \mathbb{T}^{2}, 1 \leqslant j \leqslant N_{M}$ of $\mathbb{T}^{2}$ such that the maps $y_{j}: D_{j} \rightarrow \mathbb{T}^{2}: x \mapsto(\mathrm{id}-M) x$ are bijections.

Proof. We let id $-M$ act on the cover $\mathbb{R}^{2}$. The linear map id $-M$ sends a fundamental region of $\mathbb{T}^{2}$, e.g. $[0,1)^{2}$, to a convex polytope having a non-zero volume given by $|\operatorname{det}(\mathrm{id}-M)|$. Each fixed point of $M$ on $\mathbb{T}^{2}$ is mapped by id $-M$ to an element of $\mathbb{Z}^{2}$, and the number of integer points contained in the polytope is just given by its volume. Claim (i) follows.
Let $v_{1}, v_{2} \in \mathbb{Z}^{2}$ be two different such integer points in the polytope. Now assume that there are $f_{1}, f_{2} \in[0,1)^{2}$ such that

$$
(\operatorname{id}-M)^{-1}\left(f_{1}-f_{2}\right) \equiv(\mathrm{id}-M)^{-1}\left(v_{1}-v_{2}\right) \quad\left(\bmod [0,1)^{2}\right)
$$

The right-hand side is mapped to a fixed point of $M$ on $\mathbb{T}^{2}$, implying that $f_{1}-f_{2}$ is an integer point, which is only possible if $f_{1}=f_{2}$. Therefore, $v_{1}=v_{2}$, which contradicts the assumption, and Claim (ii) follows.

For $d \in \mathbb{N}$ and every real $d \times d$ matrix $M$ we denote by $\square_{i, j}(M), 1 \leqslant i, j \leqslant d$ the submatrix arising by removing the $i$-th row and $j$-th column from $M$.

Lemma A. 2 (Determinant preserving transformation). Let $d \in \mathbb{N}$, and let $T: \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}$ and $\phi: \mathbb{R} \rightarrow \mathbb{R}$ be differentiable maps. Fix $1 \leqslant j \leqslant d$ and $\alpha \in \mathbb{R}^{d}$ and set

$$
T_{\phi}: \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}: x \mapsto\left(\alpha_{i} \phi\left(x_{j}\right) \mid 1 \leqslant i \leqslant d\right) .
$$

Then for $x \in \mathbb{R}^{d}$ it holds

$$
\operatorname{det} \mathrm{D}_{x}\left(T+T_{\phi}\right)-\operatorname{det} \mathrm{D}_{x}(T)=0
$$

## Appendix A

if and only if at least one of the conditions holds:
(i) $\sum_{i=1}^{d}(-1)^{i} \alpha_{i} \operatorname{det} \square_{i, j}\left(\mathrm{D}_{x} T\right)=0$ or
(ii) $\phi^{(1)}\left(x_{j}\right)=0$.

Proof. We develop the determinant of $\mathrm{D}_{x}\left(T+T_{\phi}\right)$ with respect to the $j$-th column. Since $T_{\phi}$ depends only on $x_{j}$ this gives

$$
\operatorname{det} \mathrm{D}_{x}\left(T+T_{\phi}\right)=(-1)^{j} \sum_{i=1}^{d}(-1)^{i} \partial_{j}\left(T+T_{\phi}\right)_{i}(x) \operatorname{det} \square_{i, j}\left(\mathrm{D}_{x} T\right)
$$

Hence, it holds

$$
\begin{aligned}
\operatorname{det} \mathrm{D}_{x}\left(T+T_{\phi}\right)-\operatorname{det}\left(\mathrm{D}_{x} T\right) & =(-1)^{j} \sum_{i=1}^{d}(-1)^{i} \operatorname{det} \square_{i, j}\left(\mathrm{D}_{x} T\right) \partial_{j}\left(T_{\phi}\right)_{i}(x) \\
& =(-1)^{j} \phi^{(1)}\left(x_{j}\right) \sum_{i=1}^{d}(-1)^{i} \alpha_{i} \operatorname{det} \square_{i, j}\left(\mathrm{D}_{x} T\right) .
\end{aligned}
$$

One deduces Claim (i) and (ii) directly from the right-hand side.
Dolgopyat pointed out that one can use some sort of bootstrapping by looking at even powers of traces to show that there are at least two non-trivial resonances. Suppose that there exists only one non-trivial $\lambda \in \operatorname{sp} \mathcal{K}_{T}$. Then we find

$$
\begin{aligned}
\operatorname{Tr} \mathcal{K}_{T}^{2} & =1+\lambda^{2}=1+\epsilon^{2} B_{M}(\psi)^{2}+O\left(2 \epsilon^{3} B_{M}(\psi)\right)+O\left(\epsilon^{4}\right) \\
& =\operatorname{Tr} \mathcal{K}_{T^{2}}=1+\epsilon B_{M^{2}}\left(\psi_{\epsilon}\right)+O\left(\epsilon^{2}\right),
\end{aligned}
$$

where $\psi_{\epsilon}:=M \psi+\psi \circ(M+\epsilon \psi)$. Then if $\psi_{0} \not \equiv 0$ would imply the existence of a further resonance (for smaller $\epsilon_{0}$ ). We could extend the reasoning by using

$$
\operatorname{Tr} \mathcal{K}_{T^{2 k}}=\operatorname{Tr} \mathcal{K}_{T}^{2^{k}} \leqslant 1+\left(\operatorname{Tr} \mathcal{K}_{T}^{2^{k-1}}-1\right)^{2}
$$

to show additional resonances for even smaller $\epsilon_{0}$.

## Appendix B

We check the expansion and contraction properties of the cones claimed in Section 2.2:

Lemma B.1. Let $C$ and $\theta$ be the constants from (2.2). Let $x \in M$ and $0<\gamma<1$ and recall the cones $C_{\gamma}^{-}(x)$ and $C_{\gamma}^{+}(x)$ defined in (2.8). Let $\alpha>0$ and $\gamma^{\prime}>0$ such that $C^{2} \theta^{\alpha} \gamma<\gamma^{\prime} \leqslant 1$. Then it holds:
(i) $\left(\mathrm{D} g_{-\alpha}\right)^{\operatorname{tr}} C_{\gamma}^{-}(x) \Subset C_{\gamma^{\prime}}^{-}\left(g_{-\alpha}(x)\right)$,
(ii) $\left(\mathrm{D} g_{\alpha}\right)^{\operatorname{tr}} C_{\gamma}^{+}(x) \Subset C_{\gamma^{\prime}}^{+}\left(g_{\alpha}(x)\right)$.

In particular, there exists $\gamma^{\prime}>0$ such that for all large enough $\alpha>0$ it holds $\gamma^{\prime}<\gamma$.

Proof. First we note that a fixed choice $\gamma^{\prime}<\gamma$ is possible for all large $\alpha>0$ because $\theta<0$. We show claim (i). Claim (ii) is shown analogously. We let $v^{-}+v^{+}+v^{0}=v \in C_{\gamma}^{-}(x)$. We estimate (assuming $\frac{1}{C}\left\|v^{0}\right\| \leqslant\left\|\left(\mathrm{D} g_{-\alpha}\right)^{\operatorname{tr}} v^{0}\right\| \leqslant$ $\left.C\left\|v^{0}\right\|\right)$

$$
\begin{aligned}
\left\|\left(\mathrm{D} g_{-\alpha}\right)^{\operatorname{tr}} v^{+}\right\|+\left\|\left(\mathrm{D} g_{-\alpha}\right)^{\operatorname{tr}} v^{0}\right\| & \leqslant C\left(\left\|v^{+}\right\|+\left\|v^{0}\right\|\right) \leqslant C \gamma\left\|v^{-}\right\| \\
& \leqslant C^{2} \theta^{\alpha} \gamma\left\|\left(\mathrm{D} g_{-\alpha}\right)^{\operatorname{tr}} v^{-}\right\|
\end{aligned}
$$

It follows that $\left(\mathrm{D} g_{-\alpha}\right)^{\operatorname{tr}} v \in C_{\gamma^{\prime}}^{-}\left(g_{-\alpha}(x)\right)$ if $\gamma^{\prime} \geqslant C^{2} \theta^{\alpha} \gamma$. Since $C_{\gamma^{\prime}}^{-}\left(g_{-\alpha}(x)\right) \Subset$ $C_{\gamma^{\prime}+\epsilon}^{-}\left(g_{-\alpha}(x)\right)$ for all $\epsilon>0$ we conclude.

Lemma B.2. Assuming the constants $C$ and $\theta$ from (2.2), let $\gamma>0, x \in M$ and suppose that $C^{2} \theta^{\alpha} \gamma<1$. Then for all $C^{2} \theta^{\alpha} \gamma<\gamma^{\prime}<1$ it holds:
(i) If $v \in C_{\gamma}^{-}(x)$ then $\frac{\left\|\left(\mathrm{D} g_{-\alpha}\right)^{\mathrm{tr}} v\right\|}{\|v\|} \geqslant C \frac{1+\gamma}{1+\gamma^{\prime}} \theta^{-\alpha}$.
(ii) If $v \in\left(\mathrm{D} g_{\alpha}\right)^{\operatorname{tr}} C_{\gamma}^{+}(x)$ then $\frac{\left\|\left(\mathrm{D} g_{-\alpha}\right)^{\operatorname{tr}} v\right\|}{\|v\|} \leqslant C \frac{1+\gamma}{1-\gamma^{\prime}} \theta^{\alpha}$.

Proof. Let $v \in T_{x}^{*} M$. We recall

$$
v=v^{-}+v^{+}+v^{0}, \quad v^{\sigma} \in E_{\sigma, x}, \sigma \in\{-,+, 0\}
$$

## Appendix $B$

If $v \in C_{\gamma}^{-}(x)$ then by (2.9), for all $\lambda \geqslant 0$ it holds

$$
\begin{aligned}
\left\|\left(\mathrm{D} g_{-\alpha}\right)^{\operatorname{tr}} v\right\| & \geqslant\left\|\left(\mathrm{D} g_{-\alpha}\right)^{\operatorname{tr}} v_{-}\right\|-\left\|\left(\mathrm{D} g_{-\alpha}\right)^{\operatorname{tr}} v_{+}\right\|-\left\|\left(\mathrm{D} g_{-\alpha}\right)^{\mathrm{tr}} v_{0}\right\| \\
& \geqslant \frac{1}{C} \theta^{-\alpha}\left\|v_{-}\right\|+C\left(\left\|v_{+}\right\|+\left\|v_{0}\right\|\right) \\
& \geqslant C \theta^{-\alpha}\left(1-\gamma^{\prime}\right)\left\|v^{-}\right\|=C \theta^{-\alpha}\left(\left(1-\gamma^{\prime}-\lambda\right)\left\|v^{-}\right\|+\lambda\left\|v^{-}\right\|\right) \\
& \geqslant C \theta^{-\alpha}\left((1+\gamma-\lambda)\left\|v^{-}\right\|+\frac{\lambda}{\gamma^{\prime}}\left(\left\|v^{0}\right\|+\left\|v^{+}\right\|\right)\right) .
\end{aligned}
$$

The choice $\lambda=\frac{1+\gamma}{\frac{1}{\gamma}+1}$ yields $\left\|\left(\mathrm{D} g_{-\alpha}\right)^{\operatorname{tr}} v\right\| \geqslant C \frac{1+\gamma}{1+\gamma^{\prime}} \theta^{-\alpha}\|v\|$.
If $v \in\left(\mathrm{D} g_{\alpha}\right)^{\operatorname{tr}} C_{\gamma}^{+}(x)$ then by (2.9), for all $\lambda \geqslant 0$ it holds

$$
\begin{aligned}
\left\|\left(\mathrm{D} g_{-\alpha}\right)^{\mathrm{tr}} v\right\| & \leqslant C \theta^{\alpha}(1+\gamma)\left\|v^{+}\right\|=C \theta^{\alpha}\left((1+\gamma+\lambda)\left\|v^{+}\right\|-\lambda\left\|v^{+}\right\|\right) \\
& \leqslant C \theta^{\alpha}\left((1+\gamma+\lambda)\left\|v^{+}\right\|-\frac{\lambda}{\gamma^{\prime}}\left(\left\|v^{0}\right\|+\left\|v^{-}\right\|\right)\right) .
\end{aligned}
$$

The choice $\lambda=\frac{1+\gamma}{\frac{\gamma^{\prime}-1}{1}}$ yields $\left\|\left(\mathrm{D} g_{-\alpha}\right)^{\operatorname{tr}} v\right\| \leqslant C \frac{1+\gamma}{1-\gamma^{\prime}} \theta^{\alpha}\|v\|$.
We let $\nabla_{z}$ be the gradient and $\nabla_{z}^{\operatorname{tr}}$ the divergence with respect to $z \in \mathbb{R}^{d}$.
Lemma B. 3 (Integration by parts (cf. [10, p.10])). Let $\mathcal{B}$ be a Banach space and let $f: \mathbb{R}^{d} \rightarrow \mathcal{B}$ be $C^{1}$ such that

$$
\|f(z)\|_{\mathcal{B}} \rightarrow 0 \quad \text { as } \quad|z| \rightarrow \infty .
$$

Let $G: \mathbb{R}^{d} \rightarrow \mathbb{R}$ be $C^{2}$ and assume that $\left|\nabla_{z} G(z)\right|>0$ for every $z \in \operatorname{supp} f$. Then it holds

$$
\int_{\mathbb{R}^{d}} e^{\mathrm{i} G(z)} f(z) \mathrm{d} z=\mathrm{i} \int_{\mathbb{R}^{d}} e^{\mathrm{i} G(z)} \nabla_{z}^{\operatorname{tr}} \frac{\nabla_{z} G(z) f(z)}{\left|\nabla_{z} G(z)\right|^{2}} \mathrm{~d} z .
$$

We understand the above transformation as integration by parts. Repeated application leads to the following iteration pattern.

Lemma B.4. Let $f(z, \eta, \xi)$ and $\nabla_{z} G(z, \eta, \xi)$ be complex and real valued functions, respectively, both $C^{r_{1}}, C^{r_{2}}, C^{r_{3}}$ in $z, \eta, \xi \in \mathbb{R}^{d}$ for some $r_{1}, r_{2}, r_{3}>0$, respectively. Let $V_{0}(z, \eta, \xi):=f(z, \eta, \xi)$ and

$$
\begin{equation*}
V_{k}(z, \eta, \xi):=\nabla_{z}^{\operatorname{tr}} \frac{\nabla_{z} G(z, \eta, \xi) V_{k-1}(z, \eta, \xi)}{\left|\nabla_{z} G(z, \eta, \xi)\right|^{2}} \text {, where } k=1, \ldots,\left\lfloor r_{1}\right\rfloor \text {. } \tag{B.1}
\end{equation*}
$$

If $\left|\nabla_{z} G(x, \eta, \xi)\right|>0$ then it holds

$$
V_{k}(z, \eta, \xi)=\left|\nabla_{z} G(z, \eta, \xi)\right|^{-k} f_{k}(z, \eta, \xi),
$$

where $f_{k}(z, \eta, \xi)$ is $C^{r_{1}-k}$ in $z, C^{r_{2}}, C^{r_{3}}$ in $\eta, \xi$, respectively and $\operatorname{supp} f=$ $\operatorname{supp} f_{k}$. Moreover, it holds for some constant $C \geqslant 1$

$$
\begin{equation*}
\left\|f_{k}\right\|_{C^{0}} \leqslant\left. C \sup _{(z, \eta, \xi) \in \operatorname{supp} f} \max _{0 \leqslant|\gamma| \leqslant k}\left|\nabla_{z} G(z, \eta, \xi)\right| \partial_{z}^{\gamma} \frac{\nabla_{z} G(z, \eta, \xi)}{\left|\nabla_{z} G(z, \eta, \xi)\right|^{2}}\right|^{k}\|f(\cdot, \eta, \xi)\|_{C^{k}} \tag{B.2}
\end{equation*}
$$

Proof. We prove this by induction. For $V_{0}=V_{0}(z, \eta, \xi)$ the hypothesis holds. We assume the hypothesis to hold for $V_{k}=V_{k}(z, \eta, \xi)$ up to some $0 \leqslant k \leqslant\left\lfloor r_{1}\right\rfloor-1$. We have therefore

$$
\begin{equation*}
V_{k+1}=\nabla_{z}^{\operatorname{tr}} \frac{\frac{\nabla_{z} G}{\left|\nabla_{z} G\right|} f_{k}}{\left|\nabla_{z} G\right|^{k+1}}=\frac{\nabla_{z}^{\operatorname{tr}}\left(\frac{\nabla_{z} G}{|\nabla z G|} f_{k}\right)}{\left|\nabla_{z} G\right|^{k+1}}-(k+1) \frac{\frac{\nabla_{z}^{\operatorname{tr}} G}{|\nabla G|^{2}} f_{k} \nabla_{z}\left|\nabla_{z} G\right|}{\left|\nabla_{z} G\right|^{k+1}} . \tag{B.3}
\end{equation*}
$$

Hence we can write $V_{k+1}=\left|\nabla_{z} G\right|^{-k-1} f_{k+1}$, where $f_{k+1}=f_{k+1}(z, \eta, \xi)$ is regular as required by the lower bound on $\left|\nabla_{z} G\right|$. In (B.3) one sees that supp $f_{k+1} \subseteq$ $\operatorname{supp} f_{k}$. From (B.3) one finds

$$
f_{k+1}=\left|\nabla_{z} G\right|^{k+1} \nabla_{z}^{\operatorname{tr}}\left(\frac{\nabla_{z} G}{\left|\nabla_{z} G\right|^{2}} \frac{f_{k}}{\left|\nabla_{z} G\right|^{k}}\right) .
$$

We recursively expand $f_{k}$ into this equation and estimate by the worst term which yields the upper bound (B.2).

A regularized version of integration by parts is used if the involved maps are only Hölder continuous. A form of Lemma B. 5 below appeared in a work of BaladiTsujii [10, p.12, Equation 3.4]. We let $\phi: \mathbb{R}^{d} \rightarrow \mathbb{R}_{\geqslant 0}$ be $C^{\infty}$, supported on the unit ball such that $\int_{\mathbb{R}^{d}} \phi(x) \mathrm{d} x=1$. For every $\epsilon>0$ we set $\phi_{\epsilon}(x)=\frac{1}{\epsilon^{d}} \phi\left(\frac{x}{\epsilon}\right)$.

Lemma B. 5 (Regularized integration by parts). Let $0<\delta<1$. Let $f: \mathbb{R}^{d} \rightarrow \mathbb{C}$ be a compactly supported $C^{\delta}$-map and let $G: \mathbb{R}^{d} \rightarrow \mathbb{R}$ be $C^{1+\delta}$ and assume that $\left|\nabla_{z} G\right|>0$ for every $z \in \operatorname{supp} f$. Set $h(z):=\frac{\nabla_{z} G(z) f(z)}{\left|\nabla_{z} G(z)\right|^{2}}$ and $h_{\epsilon}:=h * \phi_{\epsilon}$. For every $L \geqslant 1$ it holds

$$
\begin{align*}
\int_{\mathbb{R}^{d}} e^{\mathrm{i} L G(z)} f(z) \mathrm{d} z= & \frac{\mathrm{i}}{L} \int_{\mathbb{R}^{d}} e^{\mathrm{i} L G(z)} \nabla_{z}^{\operatorname{tr}} h_{\epsilon}(z) \mathrm{d} z  \tag{B.4}\\
& +\int_{\mathbb{R}^{d}} e^{\mathrm{i} L G(z)} \nabla_{z}^{\operatorname{tr}} G(z)\left(h(z)-h_{\epsilon}(z)\right) \mathrm{d} z .
\end{align*}
$$

## Appendix B

In particular, for some constant $C \geqslant 1$, it holds $\left\|\nabla_{z} h_{\epsilon}\right\|_{L_{\infty}} \leqslant C\|h\|_{C^{\delta}} \epsilon^{\delta-1}$ and $\left\|h-h_{\epsilon}\right\|_{L_{\infty}} \leqslant C\|h\|_{C^{\delta}} \epsilon^{\delta}$.

Proof. Since $G$ is $C^{1+\delta}$ and $\left|\nabla_{z} G\right|>0$, the map $h$ is $C^{\delta}$. We have $\nabla_{z}^{\operatorname{tr}} G(z) h(z)=$ $f(z)$ and we write

$$
\int_{\mathbb{R}^{d}} e^{\mathrm{i} L G(z)} f(z) \mathrm{d} z=\int_{\mathbb{R}^{d}} e^{\mathrm{i} L G(z)}\left(\nabla_{z}^{\operatorname{tr}} G(z) h_{\epsilon}(z)+\nabla_{z}^{\operatorname{tr}} G(z)\left(h(z)-h_{\epsilon}(z)\right)\right) \mathrm{d} z
$$

And since $h_{\epsilon}$ is compactly supported we have, using integration by parts,

$$
\int_{\mathbb{R}^{d}} e^{\mathrm{i} L G(z)} \nabla_{z}^{\operatorname{tr}} G(z) h_{\epsilon}(z) \mathrm{d} z=-\frac{1}{\mathrm{i} L} \int_{\mathbb{R}^{d}} e^{\mathrm{i} L G(z)} \nabla_{z}^{\operatorname{tr}} h_{\epsilon}(z) \mathrm{d} z
$$

To see the norm estimates, we have

$$
\begin{aligned}
\left|h(z)-h_{\epsilon}(z)\right| & =\left|\epsilon^{-d} \int_{\mathbb{R}^{d}}\left(h(z)-h\left(z-z^{\prime}\right)\right) \nu\left(\frac{z^{\prime}}{\epsilon}\right) \mathrm{d} z^{\prime}\right| \\
& =\left|\int_{\mathbb{R}^{d}}\left(h(z)-h\left(z-\epsilon z^{\prime}\right)\right) \nu\left(z^{\prime}\right) \mathrm{d} z^{\prime}\right| \leqslant\|h\|_{C^{\delta}} \epsilon^{\delta}
\end{aligned}
$$

Since $\operatorname{supp} h$ is compact, for every $z \in \mathbb{R}^{d}$ there exists $\bar{z} \in \mathbb{R}^{d}$ such that $h(z-\bar{z})=$ 0 . We estimate, for some constant $C \geqslant 1$, using 1-Lipschitz continuity of the norm,

$$
\begin{gathered}
\left|\int_{\mathbb{R}^{d}} h\left(z^{\prime}\right)\left(\nabla_{z} \nu\right)\left(\frac{z-z^{\prime}}{\epsilon}\right) \mathrm{d} z^{\prime}\right|=\left|\int_{\mathbb{R}^{d}}\left(h\left(z-z^{\prime}\right)-h(z-\bar{z})\right)\left(\nabla_{z} \nu\right)\left(\frac{z^{\prime}}{\epsilon}\right) \mathrm{d} z^{\prime}\right| \\
\leqslant C\|h\|_{C^{\delta}} \int_{\mathbb{R}^{d}}\left|z^{\prime}-\bar{z}\right|^{\delta}\left|\left(\nabla_{z} \nu\right)\left(\frac{z^{\prime}}{\epsilon}\right)\right| \mathrm{d} z^{\prime} \leqslant C\|h\|_{C^{\delta}} \epsilon^{\delta+d}
\end{gathered}
$$

## References

[1] A. Adam, "Generic non-trivial resonances for Anosov diffeomorphisms," Nonlinearity, vol. 30, no. 3, p. 1146, 2017.
[2] - , "Horocycle averages on closed manifolds and transfer operators," ArXiv e-prints, 2018. arXiv: 1809.04062 [math.DS].
[3] D. V. Anosov, Geodesic flows on closed Riemann manifolds with negative curvature, 1st, ser. Proceedings of the Steklov Institute of Mathematics, no. 90, 1967. American Mathematical Society, 1969.
[4] -_, "Roughness of geodesic flows on compact Riemannian manifolds of negative curvature.," Sov. Math., Dokl., vol. 3, pp. 1068-1070, 1962.
[5] V. Baladi, Dynamical Zeta Functions and Dynamical Determinants for Hyperbolic Maps: A Functional Approach. Cham: Springer, 2018, vol. 68, pp. xv +291.
[6] -, Positive transfer operators and decay of correlations. Singapore: World Scientific, 2000, pp. x +314.
[7] -_, "The quest for the ultimate anisotropic Banach space," J. Stat. Phys., vol. 166, no. 3-4, pp. 525-557, 2017.
[8] V. Baladi, M. F. Demers, and C. Liverani, "Exponential decay of correlations for finite horizon Sinai billiard flows," Inventiones mathematicae, 2017.
[9] V. Baladi and C. Liverani, "Exponential Decay of Correlations for Piecewise Cone Hyperbolic Contact Flows," Communications in Mathematical Physics, vol. 314, no. 3, pp. 689-773, 2012.
[10] V. Baladi and M. Tsujii, "Anisotropic Hölder and Sobolev spaces for hyperbolic diffeomorphisms," Ann. Inst. Fourier (Grenoble), vol. 57, pp. 127154, 2007.
[11] -, "Dynamical determinants and spectrum for hyperbolic diffeomorphisms," in Geometric and Probabilistic Structures in Dynamics, ser. Contemp. Math. Vol. 469, Amer. Math. Soc., Providence, RI, 2008, pp. 29 68.
[12] O. F. Bandtlow and F. Naud, "Lower bounds for the Ruelle spectrum of analytic expanding circle maps," Ergodic Theory and Dynamical Systems, pp. 1-22, 2017.
[13] O. F. Bandtlow, W. Just, and J. Slipantschuk", "Spectral structure of transfer operators for expanding circle maps," Annales de l'Institut Henri Poincaré (C) Non Linear Analysis, vol. 34, no. 1, pp. 31-43, 2017.
[14] P. Blanchard and E. Brüning, Mathematical Methods in Physics: Distributions, Hilbert Space Operators, Variational Methods, and Applications in Quantum Physics, 2nd ed. Birkhäuser, 2015.
[15] M. L. Blank, G. Keller, and C. Liverani, "Ruelle-Perron-Frobenius spectrum for Anosov maps," Nonlinearity, vol. 15, no. 6, pp. 1905-1973, 2002.
[16] S. Bochner, "Integration von Funktionen, deren Werte die Elemente eines Vektorraumes sind.," Fundamenta Mathematicae, vol. 20, pp. 262-276, 1933.
[17] V. I. Bogachev, Measure Theory. Springer, 2007, vol. 1.
[18] R. Bowen and B. Marcus, "Unique ergodicity for horocycle foliations," Israel Journal of Mathematics, vol. 26, no. 1, pp. 43-67, 1977.
[19] O. Butterley, "A note on operator semigroups associated to chaotic flows," Ergodic Theory and Dynamical Systems, vol. 36, no. 5, pp. 1396-1408, 2016.
[20] ——, "Dolgopyat-type bound for the resolvent," Private communication, 2018.
[21] O. Butterley and C. Liverani, "Robustly invariant sets in fiber contracting bundle flows," Journal of Modern Dynamics, vol. 7, p. 255, 2013.
[22] ——, "Smooth Anosov flows: Correlation spectra and stability," Journal of Modern Dynamics, pp. 301-322, 2007.
[23] D. Dolgopyat, "On decay of correlations in Anosov flows.," Ann. Math. (2), vol. 147, no. 2, pp. 357-390, 1998.
[24] F. Faure and C. Guillarmou, "Horocyclic invariance of Ruelle resonant states for contact Anosov flows in dimension 3," ArXiv e-prints, 2017. arXiv: 1705.07965 [math.DS].
[25] F. Faure and N. Roy, "Ruelle-Pollicott resonances for real analytic hyperbolic maps," Nonlinearity, vol. 19, no. 6, pp. 1233-1252, 2006.
[26] F. Faure and M. Tsujii, "Band structure of the Ruelle spectrum of contact Anosov flows," Comptes Rendus Mathematique, vol. 351, no. 9, pp. 385391, 2013.
[27] ——, "The semiclassical zeta function for geodesic flows on negatively curved manifolds," Inventiones mathematicae, vol. 208, no. 3, pp. 851998, 2017.
[28] L. Flaminio and G. Forni, "Invariant distributions and time averages for horocycle flows," Duke Math. J., vol. 119, no. 3, pp. 465-526, 2003.
[29] D. Fried, "Meromorphic Zeta functions for Analytic Flows," Commun. Math. Phys., vol. 174, no. 1, pp. 161-190, 1995.
[30] P. Giulietti and C. Liverani, "Parabolic dynamics and Anisotropic Banach spaces," to appear JEMS, 2018. arXiv: 1412.7181v5 [math.DS].
[31] P. Giulietti, C. Liverani, and M. Pollicott, "Anosov flows and dynamical zeta functions," Ann. Math. (2), vol. 178, no. 2, pp. 687-773, 2013.
[32] S. Gouëzel and C. Liverani, "Banach spaces adapted to Anosov systems," Ergodic Theory Dyn. Syst., vol. 26, no. 1, pp. 189-217, 2006.
[33] - , "Compact locally maximal hyperbolic sets for smooth maps: fine statistical properties," J. Differ. Geom., vol. 79, no. 3, pp. 433-477, 2008.
[34] A. Grothendieck, Produits tensoriels topologiques et espaces nucléaires, 4th ed., 16. Mem. Amer. Math. Soc., 1966.
[35] B. Hasselblatt and A. Katok, A First Course in Dynamics: with a Panorama of Recent Developments. Cambridge University Press, 2003.
[36] H. Hennion, "Sur un théorème spectral et son application aux noyaux lipchitziens," Proceedings of the American Mathematical Society, vol. 118, no. 2, pp. 627-634, 1993.
[37] E. Hille and R. S. Phillips, Functional analysis and semi-groups. 3rd printing of rev. ed. of 195\%. 1957.
[38] M. W. Hirsch and C. C. Pugh, "Smoothness of horocycle foliations.," J. Differ. Geom., vol. 10, pp. 225-238, 1975.
[39] L. Hörmander, Linear Partial Differential Operators, 4th ed., ser. Grundlehren der mathematischen Wissenschaften. Springer, 1976.
[40] S. Hurder and A. Katok, "Differentiability, rigidity and Godbillon-Vey classes for Anosov flows," Publications Mathématiques de l'IHÉS, vol. 72, pp. 5-61, 1990.
[41] A. Katok and K. Burns, "Infinitesimal Lyapunov functions, invariant cone families and stochastic properties of smooth dyanmical systems," Ergodic Theory and Dynamical Systems, vol. 14, no. 4, pp. 757-785, 1994.
[42] A. Katok and B. Hasselblatt, Handbook of Dynamical Systems. Elsevier Science, 2006, vol. 1B.
[43] ——, Introduction to the Modern Theory of Dynamical Systems, ser. Encyclopedia of mathematics and its applications. Cambridge university press, 1995.
[44] Y. Katznelson, An Introduction to Harmonic Analysis, 3 edition. Cambridge University Press, 2004.
[45] G. Keller and H. H. Rugh, "Eigenfunctions for smooth expanding circle maps," Nonlinearity, vol. 17, no. 5, pp. 1723-1730, 2004.
[46] R. N. Klaus-Jochen Engel, A Short Course on Operator Semigroups, 1st Edition., ser. Universitext. Springer, 2006.
[47] S. G. Krantz, Function Theory of Several Complex Variables, 2 edition. Amer. Math. Soc., 2001.
[48] C. Liverani, "On contact Anosov flows," Annals of Mathematics, vol. 159, no. 3, pp. 1275-1312, 2004.
[49] B. Marcus, "Ergodic Properties of Horocycle Flows for Surfaces of Negative Curvature," Annals of Mathematics, vol. 105, no. 1, pp. 81-105, 1977.
[50] - , "Unique ergodicity of the horocycle flow: variable negative curvature case.," Isr. J. Math., vol. 21, pp. 133-144, 1975.
[51] F. Naud, "Anosov diffeomorphisms with non-trivial Ruelle spectrum," Private communication, 2015.
[52] ——, "The Ruelle spectrum of generic transfer operators," Discrete Cont. Dyn. S., vol. 32, no. 7, pp. 2521-2531, 2012.
[53] B. Randol, "Small eigenvalues of the Laplace operator on compact Riemann surfaces," Bull. Amer. Math. Soc., vol. 80, no. 5, pp. 996-1000, 1974.
[54] D. Ruelle, "Zeta-Functions for Expanding Maps and Anosov Flows," Invent. Math., vol. 34, pp. 231-242, 1976.
[55] H. H. Rugh, "Generalized Fredholm determinants and Selberg zeta functions for Axiom A dynamical systems," Ergodic Theory and Dynamical Systems, vol. 16, no. 4, pp. 805-819, 1996.
[56] -_, "The correlation spectrum for hyperbolic analytic maps," Nonlinearity, vol. 5, no. 6, pp. 1237-1263, 1992.
[57] R. A. Ryan, Introduction to Tensor Products of Banach Spaces, ser. Springer Monographs in Mathematics. Springer London, 2002.
[58] D. Sarason, The $H^{p}$ Spaces of an Annulus, 56. Mem. Amer. Math. Soc., 1965.
[59] J. Slipantschuk, O. F. Bandtlow, and W. Just, "Analytic expanding circle maps with explicit spectra," Nonlinearity, vol. 26, no. 12, pp. 3231-3245, 2013.
[60] - , "Complete spectral data for analytic Anosov maps of the torus," ArXiv e-prints, 2016. arXiv: 1605.02883.
[61] - , "Complete spectral data for analytic anosov maps of the torus," Nonlinearity, vol. 30, no. 7, p. 2667, 2017.
[62] C. D. Sogge, Fourier Integrals in Classical Analysis, ser. Cambridge Tracts in Mathematics. Cambridge University Press, 1993.
[63] H. Triebel, Theory of function spaces, ser. Monographs in Mathematics. Birk-hauser, 1983.
[64] M. Tsujii, "Exponential mixing for generic volume-preserving Anosov flows in dimension three," J. Math. Soc. Japan, vol. 70, no. 2, pp. 757-821, 2018.
[65] - , "On cohomological theory of dynamical zeta functions," ArXiv eprints, 2018. arXiv: 1805.11992 [math.DS].
[66] ——, "Quasi-compactness of transfer operators for contact Anosov flows," Nonlinearity, vol. 23, no. 7, p. 1495, 2010.
[67] L.-S. Young, "What Are SRB Measures, and Which Dynamical Systems Have Them?" Journal of Statistical Physics, vol. 108, no. 5, pp. 733-754, 2002.


[^0]:    ${ }^{1}$ If $0<r<\infty$ is not an integer, $C^{r}$ means $C^{\lfloor r\rfloor}$ with all partial derivatives of order $\lfloor r\rfloor$ being ( $r-\lfloor r\rfloor$ )-Hölder continuous.

[^1]:    ${ }^{1}$ In this chapter it holds: if $r>0$, is not an integer, $C^{r}$ means $C^{\lfloor r\rfloor}$ with all partial derivatives of order $\lfloor r\rfloor$ being $(r-\lfloor r\rfloor)$-Hölder continuous.

[^2]:    ${ }^{2}$ By the support of a function $f: S \rightarrow \mathbb{C}$ we mean supp $:=\{x \in S \mid f(x) \neq 0\}$ which can be an open set in the topology of $S$.

[^3]:    ${ }^{3}$ If $\varphi \in C^{r-1}(M)$ then $\varphi_{c}:=\frac{1}{c} \int_{0}^{c} \varphi \circ g_{-\alpha} \mathrm{d} \alpha \in C_{X}^{r-1}(M)$ for all $c>0$. In the Banach spaces we construct the limit $\lim _{c \rightarrow 0} \varphi_{c}$ exists.

[^4]:    ${ }^{4}$ Note the iterated constant $C^{\check{n}}$ contrary to $C$ in [19, Assumption 3A]. This change was made to avoid a conflict in the proof of [19, Lemma 4.4], involving in there the constant $C_{6}$, and was communicated with Butterley [20].

[^5]:    ${ }^{5}$ Note that Lemma 2.4.10 and our choice of $\delta$ ensure that for any finite $b>0$ the spectral box $\Lambda_{\sigma}(b)=\left.\sigma(X+V)\right|_{W_{p}^{s, t, q}} \cap\{\Re \lambda>\delta,|\Im \lambda| \leqslant b\}$ is a finite set.

[^6]:    ${ }^{6}$ This is analogous to the decomposition in [31, Lemma 3.1]. The main difference to our decomposition is that we use a more explicit construction of the smoothing functions.

[^7]:    ${ }^{7}$ This is introduced ad hoc as it was pointed out by Colin Guillarmou and Giovanni Forni that for the weak-vanishing to imply strong vanishing is not obvious here. In some sense one would expect even a stronger statement. Namely that for every eigendistribution $\mathcal{D}$ in the expansion of Theorem 2.5.7 at least for one piece of horocycle orbit $w$ around $x \in M$ one has $\left|\gamma_{w, x}(\mathcal{D})\right|>0$.

