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Spectral Geometry

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Abstract

The goal of these lectures is to present the few fundamentals of non-commutative geometry looking around its spectral approach. Strongly motivated by physics, in particular by relativity and quantum mechanics, Chamseddine and Connes have defined an action based on spectral considerations, the so-called spectral action.

The idea is to review the necessary tools which are behind this spectral action to be able to compute it first in the case of Riemannian manifolds (Einstein–Hilbert action). Then, all primary objects defined for manifolds will be generalized to reach the level of noncommutative geometry via spectral triples, with the concrete analysis of the noncommutative torus which is a deformation of the ordinary one.

The basics of different ingredients will be presented and studied like, Dirac operators, heat equation asymptotics, zeta functions and then, how to get within the framework of operators on Hilbert spaces, the notion of noncommutative residue, Dixmier trace, pseudodifferential operators etc. These notions are appropriate in noncommutative geometry to tackle the case where the space is swapped with an algebra like for instance the noncommutative torus. Its non-compact generalization, namely the Moyal plane, is also investigated.

Motivations:

Let us first expose few motivations from physics to study noncommutative geometry which is by essence a spectral geometry. Of course, precise mathematical definitions and results will be given in the other sections.

The notion of spectrum is quite important in physics, for instance in classical mechanics, the Fourier spectrum is essential to understand vibrations or the light spectrum in electromagnetism. The notion of spectral theory is also important in functional analysis, where the spectral theorem tells us that any selfadjoint operator A can be seen as an integral over its spectral measure $A = \int_{a \in \text{Sp}(A)} a dP_a$ if $\text{Sp}(A)$ is the spectrum of A . This is of course essential in the axiomatic formulation of quantum mechanics, especially in the Heisenberg picture where the tools are the observables namely are selfadjoint operators.

But this notion is also useful in geometry. In special relativity, we consider fields $\psi(\vec{x})$ for $\vec{x} \in \mathbb{R}^4$ and the electric and magnetic fields $E, B \in \text{Function}(M = \mathbb{R}^4, \mathbb{R}^3)$. Einstein introduced in 1915 the gravitational field and the equation of motion of matter. But a problem appeared: what are the physical meaning of coordinates x^μ and equations of fields? Assume the general covariance of field equation. If $g_{\mu\nu}(x)$ or the tetradfield $e_\mu^I(x)$ is a solution (where I is a local inertial reference frame), then, for any diffeomorphism ϕ of M which is active or passive (i.e. change of coordinates), $e_\nu^I(x) = \frac{\partial x^\mu}{\partial \phi(x)^\nu} e_\mu^I(x)$ is also a solution. As a consequence, when relativity became general, the points disappeared and it remained only fields on fields in the sense that there is no fields on a given space-time. But how to practice geometry without space, given usually by a manifold M ? In this later case, the spectral approach, namely the control of eigenvalues of the scalar (or spinorial) Laplacian return important informations on M and one can even address the question if they are sufficient: can one hear the shape of M ?

There are two natural points of view on the notion of space: one is based on points (of a manifold), this is the traditional geometrical one. The other is based on algebra and this is the spectral one. So the idea is to use algebra of the dual spectral quantities.

This is of course more in the spirit of quantum mechanics but it remains to know what is a quantum geometry with bosons satisfying the Klein-Gordon equation $(\square + m^2)\psi(\vec{x}) = s_b(\vec{x})$ and fermions satisfying $(i\rlap{/}\partial - m)\psi(\vec{x}) = s_f(\vec{x})$ for sources s_b, s_f . Here $\rlap{/}\partial$ can be seen as a square root of \square and the Dirac operator will play a key role in noncommutative geometry.

In some sense, quantum forces and general relativity drive us to a spectral approach of physics, especially of space-time.

Noncommutative geometry, mainly pioneered by A. Connes (see [25, 31]), is based on a spectral triple $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ where the $*$ -algebra \mathcal{A} generalizes smooth functions on space-time M (or the coordinates) with pointwise product, \mathcal{H} generalizes the Hilbert space of above quoted spinors ψ and \mathcal{D} is a selfadjoint operator on \mathcal{H} which generalizes $\rlap{/}\partial$ via a connection on a vector bundle over M . The algebra \mathcal{A} also acts, via a representation of $*$ -algebra, on \mathcal{H} .

Noncommutative geometry treats space-time as quantum physics does for the phase-space since it gives a uncertainty principle: under a certain scale, phase-space points are indistinguishable. Below the scale Λ^{-1} , a certain renormalization is necessary. Given a geometry, the notion of action plays an essential role in physics, for instance, the Einstein-Hilbert action in gravity or the Yang-Mills-Higgs action in particle physics. So here, given the data $(\mathcal{A}, \mathcal{H}, \mathcal{D})$, the appropriate notion of action was introduced by Chamseddine and Connes [11] and defined as

$$S(\mathcal{D}, \Lambda, f) := \text{Tr} (f(\mathcal{D}/\Lambda))$$

where $\Lambda \in \mathbb{R}^+$ plays the role of a cut-off and f is a positive even function. The asymptotic series in $\Lambda \rightarrow \infty$ yields to an effective theory. For instance, this action applied to a noncommutative model of space-time $M \times F$ with a fine structure for fermions encoded in a finite geometry F gives rise from pure gravity to the standard model coupled with gravity [12, 21, 31].

The purpose of these notes is mainly to compute this spectral action on few examples like manifolds and the noncommutative torus.

In section 1, we present standard material on pseudodifferential operators over a compact Riemannian manifold. A description of the behavior of the kernel of a Ψ DO near the diagonal is given with the important example of elliptic operators. Then follows the notion of Wodzicki residue and its computation. The main point being to understand why it is a residue.

In section 2, the link with the Dixmier trace is shown. Different subspaces of compact operators are described in particular, the ideal $\mathcal{L}^{1,\infty}(\mathcal{H})$. Its definition is on purpose because in renormalization theory, one has to control the logarithmic divergency of the series $\sum_{n=1}^{\infty} n^{-1}$. We will see that this “defect” of convergence of the Riemann zeta function (in the sense that this generates a lot of complications of convergence in physics) is in fact an “advantage” because it is precisely the Dixmier trace and more generally the Wodzicki residue which are the right tools which mimics this zeta function: firstly, this controls the spectral aspects of a manifold and secondly they can be generalized to any spectral triple.

In section 3, we recall the basic definition of a Dirac (or Dirac-like) operator on a compact Riemannian manifold (M, g) endowed with a vector bundle E . An example is the (Clifford) bundle $E = \mathcal{C}\ell M$ where $\mathcal{C}\ell T_x^* M$ is the Clifford algebra for $x \in M$. This leads to the notion of spin structure, spin connection ∇^S and Dirac operator $\mathcal{D} = -ic \circ \nabla^S$ where c is the Clifford multiplication. A special focus is put on the change of metrics g under conformal transformations.

In section 4 is presented the fundamentals of heat kernel theory, namely the Green function of the heat operator $e^{t\Delta}$, $t \in \mathbb{R}^+$. In particular, its expansion as $t \rightarrow 0^+$ in terms of coefficients of the elliptic operator Δ , with a method to compute the coefficients of this expansion is explained. The idea being to replace the Laplacian Δ by \mathcal{D}^2 later on.

In section 5, a noncommutative integration theory is developed around the notion of spectral triple. This means to understand the notion of differential (or pseudodifferential) operators in this context. Within differential calculus, the link between the one-form and the fluctuations of the given \mathcal{D} is outlined.

Section 6 concerns few actions in physics, like the Einstein–Hilbert and Yang–Mills actions. The spectral action $\text{Tr}\left(f(\mathcal{D}/\Lambda)\right)$ is justified and the link between its asymptotic expansion in Λ and the heat kernel coefficients is given via the noncommutative integrals of powers of $|\mathcal{D}|$.

Section 7 gathers several results on the computation of a residue of a series of holomorphic functions, a real difficulty since one cannot commute residue and infinite sums. The notion of Diophantine condition appears and allows nevertheless this commutation for meromorphic extension of a class of zeta functions.

Section 8 is devoted to the computation of the spectral action on the noncommutative torus. After the very definitions, it is shown how to calculate with the noncommutative integral. The main technical difficulty stems from a Diophantine condition which seems necessary (but is sufficient) since any element of the smooth algebra of the torus is a series of

its generators, so the previous section is fully used. All proofs are not given, but the reader should be aware of all the main steps.

Section 9 is an approach of non-compact spectral triples. This is mandatory for physics since, a priori, the space-time is not compact. After a quick review on the difficulties which occur when $M = \mathbb{R}^d$ due to the fact that the Dirac operator has a continuous spectrum, the example of the Moyal plane is analyzed. This plane is a non-compact version of the noncom-mutative torus. Thus, no Diophantine condition appears, but the price to pay is that functional analysis is deeply used.

For each section, we suggest references since this review is by no means original.

Contents

1	Wodzicki residue and kernel near the diagonal	7
1.1	A quick overview on pseudodifferential operators	7
1.2	Case of manifolds	10
1.3	Singularities of the kernel near the diagonal	11
1.4	Wodzicki residue	15
2	Dixmier trace	20
2.1	Singular values of compact operators	20
2.2	Dixmier trace	22
3	Dirac operator	26
3.1	Definition and main properties	26
3.2	Dirac operators and change of metrics	31
4	Heat kernel expansion	34
4.1	The asymptotics of heat kernel	34
4.2	Computations of heat kernel coefficients	36
4.3	Wodzicki residue and heat expansion	37
5	Noncommutative integration	40
5.1	Notion of spectral triple	40
5.2	Notion of pseudodifferential operators	42
5.3	Zeta-functions and dimension spectrum	43
5.4	One-forms and fluctuations of \mathcal{D}	44
5.5	Tadpole	52
5.6	Commutative geometry	54
5.7	Scalar curvature	56
5.8	Tensor product of spectral triples	56
6	Spectral action	58
6.1	On the search for a good action functional	58
6.1.1	Einstein–Hilbert action	58
6.1.2	Quantum approach and spectral action	58
6.1.3	Yang–Mills action	59
6.2	Asymptotic expansion for $\Lambda \rightarrow \infty$	60
6.3	Remark on the use of Laplace transform	62
6.4	About convergence and divergence, local and global aspects of the asymptotic expansion	62
6.5	About the physical meaning of the spectral action via its asymptotics	64
7	Residues of series and integral, holomorphic continuation, etc	65
7.1	Residues of series and integral	65
7.2	Holomorphy of certain series	67
7.2.1	Proof of Lemma 7.6 for $i = 1$:	69

7.2.2	Proof of Lemma 7.6 for $i = 0$:	70
7.2.3	Proof of item (i.2) of Theorem 7.5:	71
7.2.4	Proof of item (iii) of Theorem 7.5:	71
7.2.5	Commutation between sum and residue	73
7.3	Computation of residues of zeta functions	75
7.4	Meromorphic continuation of a class of zeta functions	76
7.4.1	A family of polynomials	77
7.4.2	Residues of a class of zeta functions	77
8	The noncommutative torus	81
8.1	Definition of the nc-torus	81
8.2	Kernels and dimension spectrum	84
8.3	Noncommutative integral computations	86
8.4	The spectral action	90
8.5	Computations of f	91
8.5.1	Even dimensional case	93
8.5.2	Odd dimensional case	96
8.6	Proof of the main result	96
8.7	Beyond Diophantine equation	96
9	The non-compact case	98
9.1	The matter is not only technical	98
9.2	The Moyal product	99
9.3	The preferred unitization of the Schwartz Moyal algebra	102
9.4	The commutative case	102
9.5	The Moyal plane	103
9.5.1	The compactness condition	104
9.5.2	Spectral dimension of the Moyal planes	106

Notations:

$\mathbb{N} = \{1, 2, \dots\}$ is the set of positive integers and $\mathbb{N}_0 = \mathbb{N} \cup \{0\}$ the set of non negative integers. On \mathbb{R}^d , the volume form is $dx = dx^1 \wedge \dots \wedge dx^d$. \mathbb{S}^d is the sphere of radius one in dimension d . The induced metric:

$$d\xi = \left| \sum_{j=1}^d (-1)^{j-1} \xi_j d\xi_1 \wedge \dots \wedge \widehat{d\xi_j} \wedge \dots \wedge d\xi_d \right|$$

restricts to the volume form on \mathbb{S}^{d-1} .

M is a d -dimensional manifold with metric g .

U, V are open set either in M or in \mathbb{R}^d .

We denote by $dvol_g$ the unique volume element such that $dvol_g(\xi_1, \dots, \xi_d) = 1$ for all positively oriented g -orthonormal basis $\{\xi_1, \dots, \xi_d\}$ of $T_x M$ for $x \in M$. Thus in a local chart $\sqrt{\det g_x} |dx| = |dvol_g|$.

When $\alpha \in \mathbb{N}^d$ is a multi-index, we define

$$\partial_x^\alpha := \partial_{x_1}^{\alpha_1} \partial_{x_2}^{\alpha_2} \dots \partial_{x_d}^{\alpha_d}, \quad |\alpha| := \sum_{i=1}^d \alpha_i, \quad \alpha! := \alpha_1 \alpha_2 \dots \alpha_d.$$

For $\xi \in \mathbb{R}^d$, $|\xi| := \left(\sum_{k=1}^d |\xi_k|^2\right)^{1/2}$ is the Euclidean metric.

\mathcal{H} is a separable Hilbert space and $\mathcal{B}(\mathcal{H}), \mathcal{K}(\mathcal{H}), \mathcal{L}^p(\mathcal{H})$ denote respectively the set of bounded, compact and p -Schatten-class operators, so $\mathcal{L}^1(\mathcal{H})$ are trace-class operators.

1 Wodzicki residue and kernel near the diagonal

The aim of this section is to show that the Wodzicki's residue $WRes$ is a trace on the set $\Psi DO(M)$ of pseudodifferential operators on a compact manifold M of dimension d .

Let us first describe the steps:

- Define $WRes(P) = 2 \operatorname{Res}_{s=0} \zeta(s)$ for $P \in \Psi DO^m$ of order m and $\zeta : s \in \mathbb{C} \rightarrow \operatorname{Tr}(P\Delta^{-s})$, which is holomorphic when $\Re(s) \geq \frac{1}{2}(d+m)$.

- If $k^P(x, y)$ is the kernel of P , then its trace can be developed homogeneously as the following : $\operatorname{tr}(k^P(x, y)) = \sum_{j=-(m+d)}^0 a_j(x, x-y) - c_P(x) \log|x-y| + \dots$ where a_j is homogeneous of degree j in y and c_P is a density on M defined by $c_P(x) := \frac{1}{(2\pi)^d} \int_{\mathbb{S}^{d-1}} \operatorname{tr}(\sigma_{-d}^P(x, \xi)) d\xi$; here, σ_{-d}^P is the symbol of P of order $-d$.

The Wodzicki's residue has a simple computational form, namely $WRes P = \int_M c_P(x) |dx|$. Then, the trace property follows.

References for this section: Classical books are [101, 104]. For an orientation more in the spirit of noncommutative geometry since here we follow [88, 89] based on [3, 34], see also the excellent books [50, 84, 85, 106, 107].

1.1 A quick overview on pseudodifferential operators

In the following, $m \in \mathbb{C}$.

Definition 1.1. A symbol $\sigma(x, \xi)$ of order m is a C^∞ function: $(x, \xi) \in U \times \mathbb{R}^d \rightarrow \mathbb{C}$ satisfying

$$(i) |\partial_x^\alpha \partial_\xi^\beta \sigma(x, \xi)| \leq C_{\alpha\beta}(x) (1 + |\xi|)^{\Re(m) - |\beta|}, \quad C_{\alpha\beta} \text{ bounded on } U.$$

(ii) We suppose that $\sigma(x, \xi) \simeq \sum_{j \geq 0} \sigma_{m-j}(x, \xi)$ where σ_k is homogeneous of degree k in ξ where \simeq means a controlled asymptotic behavior

$$|\partial_x^\alpha \partial_\xi^\beta (\sigma - \sum_{j < N} \sigma_{m-j})(x, \xi)| \leq C_{N\alpha\beta}(x) |\xi|^{\Re(m) - N - |\beta|} \text{ for } |\xi| \geq 1 \text{ with } C_{N\alpha\beta} \text{ bounded on } U.$$

The set of symbols of order m is denoted by $S^m(U \times \mathbb{R}^d)$.

A function $a \in C^\infty(U \times U \times \mathbb{R}^d)$ is an amplitude of order m , if for any compact $K \subset U$ and any $\alpha, \beta, \gamma \in \mathbb{N}^d$ there exists a constant $C_{K\alpha\beta\gamma}$ such that

$$|\partial_x^\alpha \partial_y^\gamma \partial_\xi^\beta a(x, y, \xi)| \leq C_{K\alpha\beta\gamma} (1 + |\xi|)^{\Re(m) - |\beta|}, \quad \forall x, y \in K, \xi \in \mathbb{R}^d.$$

The set of amplitudes is written $A^m(U)$.

For $\sigma \in S^m(U \times \mathbb{R}^d)$, we get a continuous operator $\sigma(\cdot, D) : u \in C_c^\infty(U) \rightarrow C^\infty(U)$ given by

$$\sigma(\cdot, D)(u)(x) := \sigma(x, D)(u) := \frac{1}{(2\pi)^d} \int_{\mathbb{R}^d} \sigma(x, \xi) \hat{u}(\xi) e^{ix \cdot \xi} d\xi \quad (1)$$

where $\hat{\cdot}$ means the Fourier transform. This operator $\sigma(\cdot, D)$ will be also denoted by $Op(\sigma)$. For instance,

$$\text{if } \sigma(x, \xi) = \sum_{\alpha} a_{\alpha}(x) \xi^{\alpha}, \text{ then } \sigma(x, D) = \sum_{\alpha} a_{\alpha}(x) D_x^{\alpha} \text{ with } D_x := -i\partial_x.$$

Remark that, by transposition, there is a natural extension of $\sigma(\cdot, D)$ from the set $\mathcal{D}'_c(U)$ of distributions with compact support in U to the set of distributions $\mathcal{D}'(U)$.

By definition, the leading term for $|\alpha| = m$ is the *principal symbol* and the *Schwartz kernel* of $\sigma(x, D)$ is defined by

$$k^{\sigma(x, D)}(x, y) := \frac{1}{(2\pi)^d} \int_{\mathbb{R}^d} \sigma(x, \xi) e^{i(x-y)\cdot\xi} d\xi = \check{\sigma}_{\xi \rightarrow y}(x, x-y) \quad (2)$$

where $\check{\cdot}$ is the Fourier inverse in variable ξ . Similarly, if the kernel of the operator $Op(a)$ associated to the amplitude a is

$$k^a(x, y) := \frac{1}{(2\pi)^d} \int_{\mathbb{R}^d} a(x, y, \xi) e^{i(x-y)\cdot\xi} d\xi. \quad (3)$$

Definition 1.2. $P : C_c^{\infty}(U) \rightarrow C^{\infty}(U)$ (or $\mathcal{D}'(U)$) is said to be *smoothing* if its kernel is in $C^{\infty}(U \times U)$ and $\Psi DO^{-\infty}(U)$ will denote the set of smoothing operators.

For $m \in \mathbb{C}$, the set $\Psi DO^m(U)$ of pseudodifferential operators of order m will be the set of P such that

$$P : C_c^{\infty}(U) \rightarrow C^{\infty}(U), Pu(x) = (\sigma(x, D) + R)(u) \text{ where } \sigma \in S^m(U \times \mathbb{R}^d), R \in \Psi DO^{-\infty}.$$

σ is called the *symbol* of P .

Remark 1.3. It is important to quote that a smoothing operator is a pseudodifferential operator whose amplitude is in $A^m(U)$ for all $m \in \mathbb{R}$: by (3), $a(x, y, \xi) := e^{-i(x-y)\cdot\xi} k(x, y) \phi(\xi)$ where the function $\phi \in C_c^{\infty}(\mathbb{R}^d)$ satisfies $\int_{\mathbb{R}^d} \phi(\xi) d\xi = (2\pi)^d$.

Clearly, the main obstruction to smoothness is on the diagonal since

Lemma 1.4. $k^{\sigma(x, D)}$ is C^{∞} outside the diagonal.

Proof. $\check{\sigma}$ is smooth since it is given for $y \neq 0$ by the oscillatory integral

$$\int_{\mathbb{R}^d} \sigma(x, \xi) e^{iy\cdot\xi} d\xi = \int_{\mathbb{R}^d} (P_y^k \sigma) e^{iy\cdot\xi} d\xi$$

where k is an integer such that $k > \Re(m) + n$ and $P_y = P(y, D_{\xi})$ is chosen with $P_y(e^{iy\cdot\xi}) = e^{iy\cdot\xi}$; for instance $P_y = \frac{1}{|y|^2} \sum_j y_j \frac{\partial}{\partial \xi_j}$. The last integral is absolutely converging. \square

Few remarks on the duality between symbols and pseudodifferential operators:

$$\sigma(x, \xi) \in S^m(U \times \mathbb{R}^d) \longleftrightarrow k_{\sigma}(x, y) \in C_c^{\infty}(U \times U \times \mathbb{R}^d) \longleftrightarrow A = Op(\sigma) \in \Psi DO^m$$

where we used the following definition

$$\sigma^A(x, \xi) := e^{-ix\cdot\xi} A(x \rightarrow e^{ix\cdot\xi}).$$

Moreover,

$$\begin{aligned}\sigma^A &\simeq \sum_{\alpha} \frac{(-i)^{\alpha}}{\alpha!} \partial_{\xi}^{\alpha} \partial_y^{\alpha} k_{\sigma}^A(x, y, \xi)|_{y=x}, \\ k_{\sigma}^A(x, y) &=: \frac{1}{(2\pi)^d} \int_{\mathbb{R}^d} e^{i(x-y)\cdot\xi} k^A(x, y, \xi) d\xi,\end{aligned}$$

where $k^A(x, y, \xi)$ is the amplitude of $k_{\sigma}^A(x, y)$. Actually, $\sigma^A(x, \xi) = e^{iD_{\xi} D_y} k^A(x, y, \xi)|_{y=x}$ and $e^{iD_{\xi} D_y} = 1 + iD_{\xi} D_y - \frac{1}{2}(D_{\xi} D_y)^2 + \dots$. Thus $A = Op(\sigma^A) + R$ where R is a regularizing operator on U .

A point of interest is that differential operators are local: if $f = 0$ on U^c (complementary set of U) then $Pf = 0$ on U^c . While pseudodifferential operators are pseudo-local: Pf is smooth on U when f is smooth.

There are two fundamental points about ΨDO 's: they form an algebra and this notion is stable by diffeomorphism justifying its extension to manifolds and then to bundles:

Theorem 1.5. (i) If $P_1 \in \Psi DO^{m_1}$ and $P_2 \in \Psi DO^{m_2}$, then $P_1 P_2 \in \Psi DO^{m_1+m_2}$ with symbol

$$\sigma^{P_1 P_2}(x, \xi) \simeq \sum_{\alpha \in \mathbb{N}^d} \frac{(-i)^{\alpha}}{\alpha!} \partial_{\xi}^{\alpha} \sigma^{P_1}(x, \xi) \partial_x^{\alpha} \sigma^{P_2}(x, \xi).$$

The principal symbol of $P_1 P_2$ is

$$\sigma_{m_1+m_2}^{P_1 P_2}(x, \xi) = \sigma_{m_1}^{P_1}(x, \xi) \sigma_{m_2}^{P_2}(x, \xi).$$

(ii) Let $P \in \Psi DO^m(U)$ and $\phi \in \text{Diff}(U, V)$ where V is another open set of \mathbb{R}^d . The operator $\phi_* P : f \in C^{\infty}(V) \rightarrow P(f \circ \phi) \circ \phi^{-1}$ satisfies $\phi_* P \in \Psi DO^m(V)$ and its symbol is

$$\sigma^{\phi_* P}(x, \xi) = \sigma_m^P(\phi^{-1}(x), (d\phi)^t \xi) + \sum_{|\alpha| > 0} \frac{(-i)^{\alpha}}{\alpha!} \phi_{\alpha}(x, \xi) \partial_{\xi}^{\alpha} \sigma^P(\phi^{-1}(x), (d\phi)^t \xi)$$

where ϕ_{α} is a polynomial of degree α in ξ . Moreover, its principal symbol is

$$\sigma_m^{\phi_* P}(x, \xi) = \sigma_m^P(\phi^{-1}(x), (d\phi)^t \xi).$$

In other terms, the principal symbol is covariant by diffeomorphism: $\sigma^{\phi_* P}_m = \phi_* \sigma_m^P$.

While the proof of formal expressions is a direct computation, the asymptotic behavior requires some care, see [101, 104].

An interesting remark is in order: $\sigma^P(x, \xi) = e^{-ix \cdot \xi} P(x \rightarrow e^{ix \cdot \xi})$, thus the dilation $\xi \rightarrow t\xi$ with $t > 0$ gives

$$t^{-m} e^{-itx \cdot \xi} P e^{itx \cdot \xi} = t^{-m} \sigma^P(x, t\xi) \simeq t^{-m} \sum_{j \geq 0} \sigma_{m-j}^P(x, t\xi) = \sigma_m^P(x, \xi) + o(t^{-1}).$$

Thus, if $P \in \Psi DO^m(U)$ with $m \geq 0$,

$$\sigma_m^P(x, \xi) = \lim_{t \rightarrow \infty} t^{-m} e^{-ith(x)} P e^{ith(x)}, \text{ where } h \in C^{\infty}(U) \text{ is (almost) defined by } dh(x) = \xi.$$

1.2 Case of manifolds

Let M be a (compact) Riemannian manifold of dimension d . Thanks to Theorem 1.5, the following makes sense:

Definition 1.6. $\Psi DO^m(M)$ is defined as the set of operators $P : C_c^\infty(M) \rightarrow C^\infty(M)$ such that

- (i) the kernel $k^P \in C^\infty(M \times M)$ off the diagonal,
 - (ii) The map $f \in C_c^\infty(\phi(U)) \rightarrow P(f \circ \phi) \circ \phi^{-1} \in C^\infty(\phi(U))$ is in $\Psi DO^m(\phi(U))$
- for every coordinate chart $(U, \phi : U \rightarrow \mathbb{R}^d)$.

Of course, this can be generalized:

Definition 1.7. Given a vector bundle E over M , a linear map $P : \Gamma_c^\infty(M, E) \rightarrow \Gamma^\infty(M, E)$ is in $\Psi DO^m(M, E)$ when k^P is smooth off the diagonal, and local expressions are ΨDO 's with matrix-valued symbols.

The covariance formula implies that σ_m^P is independent of the chosen local chart so is globally defined on the bundle $T^*M \rightarrow M$ and σ_m^P is defined for every $P \in \Psi DO^m$ using overlapping charts and patching with partition of unity.

An important class of pseudodifferential operators are those which are invertible modulo regularizing ones:

Definition 1.8. $P \in \Psi DO^m(M, E)$ is elliptic if $\sigma_m^P(x, \xi)$ is invertible for all $\xi \in TM_x^*$, $\xi \neq 0$.

This means that $|\sigma^P(x, \xi)| \geq c_1(x)|\xi|^m$ for $|\xi| \geq c_2(x)$, $x \in U$ where c_1, c_2 are strictly positive continuous functions on U .

This also means that there exists a *parametrix*:

Lemma 1.9. The following are equivalent:

- (i) $Op(\sigma) \in \Psi DO^m(U)$ is elliptic.
- (ii) There exist $\sigma' \in S^{-m}(U \times \mathbb{R}^d)$ such that $\sigma \circ \sigma' = 1$ or $\sigma' \circ \sigma = 1$.
- (iii) $Op(\sigma) Op(\sigma') = Op(\sigma') Op(\sigma) = 1$ modulo $\Psi DO^{-\infty}(U)$.

Thus $Op(\sigma') \in \Psi DO^{-m}(U)$ is also elliptic.

At this point, it is useful to remark that any $P \in \Psi DO(M, E)$ can be extended to a bounded operator on $L^2(M, E)$ when $\Re(m) \leq 0$. Of course, this needs an existing scalar product for given metrics on M and E .

Theorem 1.10. When $P \in \Psi DO^{-m}(M, E)$ is elliptic with $\Re(m) > 0$, its spectrum is discrete when M is compact.

Proof. We need to get the result first for an open set U , for a manifold M and then for a bundle E over M .

For any $s \in \mathbb{R}$, the usual Sobolev spaces $H^s(\mathbb{R}^d)$ (with $H^0(\mathbb{R}^d) = L^2(\mathbb{R}^d)$) and $H_c^s(U)$ (defined as the union of all $H^s(K)$ over compact subsets $K \subset U$) or $H_{loc}^s(U)$ (defined as the set of distributions $u \in \mathcal{D}'(U)$ such that $\phi u \in H^s(\mathbb{R}^d)$ for all $\phi \in C_c^\infty(U)$) can be extended for any manifold M to the Sobolev spaces $H_c^s(M)$ (obvious definition) and $H_{loc}^s(M)$: if $(U, \phi : U \rightarrow \mathbb{R}^d)$ is a local chart and $\chi \in C_c^\infty(\phi(U))$, we say that a distribution $u \in \mathcal{D}'(M)$

is in $H_{loc}^s(M)$ when $(\phi^{-1})^*(\xi u) \in H^s(\mathbb{R}^d)$. When M is compact, $H_{loc}^s(M) = H_c^s(M)$ (thus denoted $H^s(M)$). Using Rellich's theorem, the inclusion $H_c^s(U) \hookrightarrow H_c^t(U)$ for $s < t$ is compact. Since $P : H_c^s(M) \rightarrow H_{loc}^{s-\Re(m)}(M)$ is a continuous linear map for a (non-necessarily compact) manifold M , both results yield that P is compact. Finally, the extended operator on a bundle is $P : L^2(M, E) \rightarrow H^{-\Re(m)}(M, E) \hookrightarrow L^2(M, E)$ where the second map is the continuous inclusion, so P being compact as an L^2 operator has a discrete spectrum. \square

We rephrase a previous remark (see [4, Proposition 2.1]):

Let E be a vector bundle of rank r over M . If $P \in \Psi DO^{-m}(M, E)$, then for any couple of sections $s \in \Gamma^\infty(M, E)$, $t^* \in \Gamma^\infty(M, E^*)$, the operator $f \in C^\infty(M) \rightarrow \langle t^*, P(fs) \rangle \in C^\infty(M)$ is in $\Psi DO^m(M)$. This means that in a local chart (U, ϕ) , these operators are $r \times r$ matrices of pseudodifferential operators of order $-m$. The total symbol is in $C^\infty(T^*U) \otimes \text{End}(E)$ with $\text{End}(E) \simeq M_r(\mathbb{C})$. The principal symbol can be globally defined: $\sigma_{-m}^P(x, \xi) : E_x \rightarrow E_x$ for $x \in M$ and $\xi \in T_x^*M$, can be seen as a smooth homomorphism homogeneous of degree $-m$ on all fibers of T^*M . Moreover, we get the simple formula which could be seen as a definition of the principal symbol

$$\sigma_{-m}^P(x, \xi) = \lim_{t \rightarrow \infty} t^{-m} \left(e^{-ith} \cdot P \cdot e^{ith} \right)(x) \text{ for } x \in M, \xi \in T_x^*M \quad (4)$$

where $h \in C^\infty(M)$ is such that $d_x h = \xi$.

1.3 Singularities of the kernel near the diagonal

The question to be solved is to define a homogeneous distribution which is an extension on \mathbb{R}^d of a given homogeneous symbol on $\mathbb{R}^d \setminus \{0\}$. Such extension is a regularization used for instance by Epstein–Glaser in quantum field theory.

The Schwartz space on \mathbb{R}^d is denoted by \mathcal{S} and the space of tempered distributions by \mathcal{S}' .

Definition 1.11. For $f_\lambda(\xi) := f(\lambda\xi)$, $\lambda \in \mathbb{R}_+^*$, define $\tau \in \mathcal{S}' \rightarrow \tau_\lambda$ by $\langle \tau_\lambda, f \rangle := \lambda^{-d} \langle \tau, f_{\lambda^{-1}} \rangle$ for all $f \in \mathcal{S}$.

A distribution $\tau \in \mathcal{S}'$ is homogeneous of order $m \in \mathbb{C}$ when $\tau_\lambda = \lambda^m \tau$.

Proposition 1.12. Let $\sigma \in C^\infty(\mathbb{R}^d \setminus \{0\})$ be a homogeneous symbol of order $k \in \mathbb{Z}$.

(i) If $k > -d$, then σ defines a homogeneous distribution.

(ii) If $k = -d$, there exists a unique obstruction to the extension of σ given by

$$c_\sigma = \int_{\mathbb{S}^{d-1}} \sigma(\xi) d\xi,$$

namely, one can at best extend σ in $\tau \in \mathcal{S}'$ such that

$$\tau_\lambda = \lambda^{-d} \left(\tau + c_\sigma \log(\lambda) \delta_0 \right). \quad (5)$$

Proof. (i) For $k > -d$, σ is integrable near zero, increases slowly at ∞ , so defines by extension a unique distribution $\tau \in \mathcal{S}'$ which will be homogeneous of order k .

(ii) Assume $k = -d$. Then σ extends to a continuous linear form $L_\sigma(f) := \int_{\mathbb{R}^d} f(\xi) \sigma(\xi) d\xi$ on $\mathcal{S}_0 := \{f \in \mathcal{S} \mid f(0) = 0\}$. By Hahn–Banach theorem, L_σ extends to \mathcal{S}' and $L_\sigma \in E$ where $E := \{\tau \in \mathcal{S}' \mid \tau|_{\mathcal{S}_0} = L_\sigma\}$ is given by the direction δ_0 .

This affine space E is stable by the endomorphism $\tau \rightarrow \lambda^d \tau_\lambda$: actually if $f \in \mathcal{S}_0$, $f_{\lambda^{-1}} \in \mathcal{S}_0$ and

$$\lambda^d \langle \tau_\lambda, f \rangle = \langle \tau, f_{\lambda^{-1}} \rangle = L_\sigma(f_{\lambda^{-1}}) = \int_{\mathbb{R}^d} f(\lambda^{-1}\xi) \sigma(\xi) d\xi = \int_{\mathbb{R}^d} f(\xi) \sigma(\xi) d\xi = L_\sigma(f),$$

thus $\lambda^d \tau_\lambda = L_\sigma$ on \mathcal{S}_0 .

Moreover, $\lambda^d (\delta_0)_\lambda = \delta_0$; thus there exists $c(\lambda) \in \mathbb{C}$ such that

$$\tau_\lambda = \lambda^{-d} \tau + c(\lambda) \lambda^{-d} \delta_0 \quad (6)$$

for all $\tau \in E$. The computation of $c(\lambda)$ for a specific example in E gives $c(\lambda) = c_\sigma \log(\lambda)$: for instance, choose $g \in C_c^\infty([0, \infty])$ which is 1 near 0 and define $\tau \in \mathcal{S}'$ by

$$\langle \tau, f \rangle := L_\sigma(f - f(0)g(|\cdot|)) = \int_{\mathbb{R}^d} (f(\xi) - f(0)g(|\xi|)) \sigma(\xi) d\xi, \quad \forall f \in \mathcal{S}.$$

Thus if $f(0) = 1$, we get $c(\lambda) \lambda^{-d} \langle \delta_0, f \rangle = c(\lambda) \lambda^{-d}$, so by (6)

$$\begin{aligned} c(\lambda) \lambda^{-d} &= \langle \tau, f_{\lambda^{-1}} \rangle - \lambda^{-d} \langle \tau, f \rangle \\ &= \int_{\mathbb{R}^d} (f(\lambda^{-1}\xi) - g(|\xi|)) \sigma(\xi) d\xi - \lambda^{-d} \int_{\mathbb{R}^d} (f(\xi) - g(|\xi|)) \sigma(\xi) d\xi \\ &= -\lambda^{-d} \int_{\mathbb{R}^d} (g(\lambda|\xi|) - g(|\xi|)) \sigma(\xi) d\xi = -\lambda^{-d} c_\sigma \int_0^\infty (g(\lambda|\xi|) - g(|\xi|)) \frac{d|\xi|}{|\xi|} \end{aligned}$$

with $c_\sigma := \int_{\mathbb{S}^{d-1}} \sigma(\xi) d^{d-1}\xi$. Since

$$\lambda \frac{d}{d\lambda} \int_0^\infty (g(\lambda|\xi|) - g(|\xi|)) \frac{d|\xi|}{|\xi|} = \lambda \int_0^\infty g'(\lambda|\xi|) d|\xi| = -g(0) = -1,$$

we get $c(\lambda) = c_\sigma \log(\lambda)$. Thus, when $c_\sigma = 0$, every element of E is a homogeneous distribution on \mathbb{R}^d which extends the symbol σ .

Conversely, let $\tau \in \mathcal{S}'$ be a homogeneous distribution extending σ and let $\tilde{\tau} \in E$. Since $\tau - \tilde{\tau}$ is supported at the origin, we can write $\tau = \tilde{\tau} + \sum_{|\alpha| \leq N} a_\alpha \partial^\alpha \delta_0$ where $a_\alpha \in \mathbb{C}$ and

$$0 = \tau_\lambda - \lambda^{-d} \tau = c_\sigma \lambda^{-d} \log(\lambda) \delta_0 + \sum_{1 \leq |\alpha| \leq N} a_\alpha \lambda^{-d} (\lambda^{|\alpha|} - 1) \partial^\alpha \delta_0.$$

The linear independence of $(\partial^\alpha \delta_0)$ gives $a_\alpha = 0$, $\forall a_\alpha$. So $c_\sigma = 0$ and $\tau \in E$. The condition $c_\sigma = 0$ is so necessary and sufficient to extend σ in a homogeneous distribution. And in the general case, one can at best extend it in a distribution satisfying (5), but it is only possible with elements of E . \square

In the following result, we are interested by the behavior near the diagonal of the kernel k^P for $P \in \Psi DO$. For any $\tau \in \mathcal{S}'$, we choose the decomposition as $\tau = \phi \circ \tau + (1 - \phi) \circ \tau$ where $\phi \in C_c^\infty(\mathbb{R}^d)$ and $\phi = 1$ near 0. We can look at the infrared behavior of τ near the origin and its ultraviolet behavior near infinity. Remark first that, since $\phi \circ \tau$ has a compact support, $(\phi \circ \tau)^\vee \in \mathcal{S}'$, so the regularity of τ^\vee depends only of its ultraviolet part $((1 - \phi) \circ \tau)^\vee$.

Proposition 1.13. *Let $P \in \Psi DO^m(U)$, $m \in \mathbb{Z}$. Then, in local form near the diagonal,*

$$k^P(x, y) = \sum_{-(m+d) \leq j \leq 0} a_j(x, x-y) - c_P(x) \log|x-y| + \mathcal{O}(1)$$

where $a_j(x, y) \in C^\infty(U \times U \setminus \{x\})$ is homogeneous of order j in y and $c_P(x) \in C^\infty(U)$ is given by

$$c_P(x) = \frac{1}{(2\pi)^d} \int_{\mathbb{S}^{d-1}} \sigma_{-d}^P(x, \xi) d\xi. \quad (7)$$

Proof. We know that $\sigma^P(x, \xi) \simeq \sum_{j \leq m} \sigma_j^P(x, \xi)$ and by (2), $k^P(x, y) = \widetilde{\sigma_{\xi \rightarrow y}^P}(x, x-y)$ so we need to control $\widetilde{\sigma_{\xi \rightarrow y}^P}(x, x-y)$ when $y \rightarrow 0$.

Assume first that $\sigma^P(x, \xi)$ is independent of x :

For $-d < j \leq m$, $\sigma_j(\xi)$ extends to $\tau_j \in \mathcal{S}'$. For $j > -d$, this extension is homogeneous (of degree j) and unique.

For $j = -d$, we may assume that τ_{-d} satisfies (5). Thus $\tau := \sigma^P - \sum_{j=-d}^m \tau_j \in \mathcal{S}'$ behaves in the ultraviolet as a integrable symbol. In particular $\widetilde{\tau}$ is continuous near 0 and we get

$$\widetilde{\sigma^P}(y) = \sum_{j=-d}^m \widetilde{\tau}_j(y) + \mathcal{O}(1). \quad (8)$$

Note that the inverse Fourier transform of the infrared part of τ_j is in $C^\infty(\mathbb{R}^d)$ while those of its ultraviolet part is in $C^\infty(\mathbb{R}^d \setminus \{0\})$, so $\widetilde{\tau}_j$ is smooth near 0.

Moreover, for $j > -d$, $\widetilde{\tau}_j$ is homogeneous of degree $-(d+j)$ while for $j = -d$,

$$\widetilde{\tau}_{-d}(\lambda y) = \lambda^{-d} [(\tau_{-d})_{\lambda^{-1}}]^\vee(y) = [\tau_{-d} - c_{\sigma_{-d}} \log(\lambda) \delta_0]^\vee(y) = \widetilde{\tau}_{-d}(y) - \frac{1}{(2\pi)^d} c_{\sigma_{-d}} \log \lambda.$$

For $\lambda = |y|^{-1}$, we get

$$\widetilde{\tau}_{-d}\left(\frac{y}{|y|}\right) = \widetilde{\tau}_{-d}\left(\frac{y}{|y|}\right) - \frac{1}{(2\pi)^d} c_{\sigma_{-d}} \log |y|.$$

Summation over j in (8) yields the result.

Assume now that $\sigma^P(x, \xi)$ is independent of x :

We do the same with families $\{\tau_x\}_{x \in U}$ and $\{\tau_{j,x}\}_{x \in U}$. Their ultraviolet behaviors are those of smooth symbols on $U \times \mathbb{R}^d$, so given by smooth functions on $U \times \mathbb{R}^d \setminus \{0\}$ and for τ_x by a continuous function on $U \times \mathbb{R}^d$. For the infrared part, we get smooth maps from U to $\mathcal{E}(\mathbb{R}^d)'$ (distributions with compact support), thus applying inverse Fourier transform, we end up with smooth functions on $U \times \mathbb{R}^d$. Actually, for $\tau_{j,x}$ with $j > -d$, this follows from the fact that it is the extension of $\sigma_j(x, \cdot)$ which is integrable near the origin: let $f \in \mathcal{S}$,

$$\langle \phi \circ \tau_{j,x}, f \rangle = \langle \tau_{j,x}, \phi \circ f \rangle = \int_{\mathbb{R}^d} \phi(\xi) f(\xi) \sigma_j(x, \xi) d\xi.$$

While for $j = -d$,

$$\langle \phi \circ \tau_{-d,x}, f \rangle = \langle \tau_{-d,x}, \phi \circ f \rangle = \int_{\mathbb{R}^d} \phi(\xi) (f(\xi) - f(0)) \sigma_{-d}(x, \xi) d\xi,$$

and the map $x \rightarrow \phi \circ \tau_{-d,x}$ is smooth from U to $\mathcal{E}(\mathbb{R}^d)'$. In conclusion,

$$\check{\sigma}_{\xi \rightarrow y}(x, y) = \sum_{-(m+d) \leq j \leq 0} a_j(x, y) - c_P(x) \log |y| + R(x, y)$$

where $a_j(x, y)$ is a smooth function on $U \times \mathbb{R}^d \setminus \{0\}$, is homogeneous of degree j in y , c_P is given by (7) and $R(x, y)$ is a function, continuous on $U \times \mathbb{R}^d$. So we get the claimed asymptotic behavior. \square

Theorem 1.14. *Let $P \in \Psi DO^m(M, E)$ with $m \in \mathbb{Z}$. Then, for any trivializing local coordinates*

$$\text{tr}(k^P(x, y)) = \sum_{j=-(m+d)}^0 a_j(x, x-y) - c_P(x) \log |x-y| + \mathcal{O}(1),$$

where a_j is homogeneous of degree j in y , c_P is intrinsically locally defined by

$$c_P(x) := \frac{1}{(2\pi)^d} \int_{\mathbb{S}^{d-1}} \text{tr}(\sigma_{-d}^P(x, \xi)) d\xi. \quad (9)$$

Moreover, $c_P(x)|dx|$ is a 1-density over M which is functorial with respect to diffeomorphisms ϕ :

$$c_{\phi_* P}(x) = \phi_* (c_P(x)). \quad (10)$$

Proof. The asymptotic behavior follows from Proposition 1.13 but we first have to understand why c_P is well defined:

Assume first that E is a trivial line bundle and P is a scalar ΨDO .

Define a change of coordinates by $y := \phi^{-1}(x)$. Thus $k^P(x, x') \xrightarrow{\phi_*} k^{\phi_* P}(y, y')$ with

$$\begin{aligned} k^{\phi_* P}(y, y') &= |J_\phi(y')| k^P(\phi(y), \phi(y')) \\ &= \sum_{j=-(m+d)}^0 |J_\phi(y')| \left[a_j(\phi(y), \phi(y) - \phi(y')) - c_P(\phi(y)) \log |\phi(y) - \phi(y')| \right] + \mathcal{O}(1). \end{aligned}$$

A Taylor expansion around $(\phi(y), \phi'(y) \cdot (y - y'))$ of a_j gives

$$a_j(\phi(y), \phi(y) - \phi(y')) \simeq |y - y'|^j a_j(\phi(y), \phi'(y) \cdot \frac{y - y'}{|y - y'|}) + \dots,$$

since $a_j(\phi(y), \cdot)$ is smooth outside 0, so we get only homogeneous and continuous terms. Moreover the only contribution to the log-term is

$$|J_\phi(y')| c_P(\phi(y)) \log |\phi(y) - \phi(y')| \simeq |J_\phi(y)| c_P(\phi(y)) \log |\phi(y) - \phi(y')| + \mathcal{O}(1)$$

and we get

$$c_{\phi_* P}(y) = |J_\phi(y)| c_P(\phi(y)).$$

In particular $c_P(x) |dx^1 \wedge \dots \wedge dx^d|$ can be globally defined on M as a 1-density. (Recall that a α -density on a vector space E of dimension n is any application $f : \Lambda^n E \rightarrow \mathbb{R}$ such

that for any $\lambda \in \mathbb{R}$, $f(\lambda x) = |\lambda|^\alpha f(x)$ and the set of these densities is denoted $|\Lambda|^\alpha E^*$; this is generalized to a vector bundle E over M where each fiber is $|\Lambda|^\alpha E_x^*$. The interest of the bundle of 1-densities is to give a class of objects directly integrable on M . In particular, we get here something intrinsically defined, even when the manifold is not oriented).

General case:

P acts on section of a bundle. By a change of trivialization, the action of P is conjugate on each fiber by a smooth matrix-valued map $A(x)$, so $k^P(x, x') \rightarrow A(x)^{-1}k^P(x, x')A(x')$. We are looking for the logarithmic term: only the 0-order term in $A(x')$ will contribute and $\text{tr}(A(x)^{-1}k^P(x, x')A(x'))$ has the same logarithmic singularity than the similar term $\text{tr}(A(x)^{-1}k^P(x, x')A(x)) = \text{tr}(k^P(x, x'))$ near the diagonal. Thus $c_P(x)$ is independent of a chosen trivialization.

Similarly, if P is not a scalar but $\text{End}(E)$ -valued, the above proof can be generalized (the space of $C^\infty((M, |\Lambda|(M) \otimes \text{End}(E))$ of $\text{End}(E)$ -valued densities is a sheaf). \square

Remark that, when M is Riemannian with metric g and $d_g(x, y)$ is the geodesic distance, then

$$\text{tr}(k^P(x, y)) = \sum_{j=-(m+d)}^0 a_j(x, x-y) - c_P(x) \log(d_g(x, y)) + \mathcal{O}(1),$$

since there exists $c > 0$ such that $c^{-1}|x-y| \leq d_g(x, y) \leq c|x-y|$.

1.4 Wodzicki residue

The claim is that $\int_M c_P(x)|dx|$ is a residue.

For this, we embed everything in \mathbb{C} . In the same spirit as in Proposition 1.12, one obtains the following

Lemma 1.15. *Every $\sigma \in C^\infty(\mathbb{R}^d \setminus \{0\})$ which is homogeneous of degree $m \in \mathbb{C} \setminus \mathbb{Z}$ can be uniquely extended to a homogeneous distribution.*

Definition 1.16. *Let U be an open set in \mathbb{R}^d and Ω be a domain in \mathbb{C} .*

A map $\sigma : \Omega \rightarrow S^m(U \times \mathbb{R}^d)$ is said to be holomorphic when

the map: $z \in \Omega \rightarrow \sigma(z)(x, \xi)$ is analytic for all $x \in U$, $\xi \in \mathbb{R}^d$,

the order $m(z)$ of $\sigma(z)$ is analytic on Ω ,

the two bounds of Definition 1.1 (i) and (ii) of the asymptotics $\sigma(z) \simeq \sum_j \sigma_{m(z)-j}(z)$ are locally uniform in z .

This hypothesis is sufficient to get:

The map: $z \rightarrow \sigma_{m(z)-j}(z)$ is holomorphic from Ω to $C^\infty(U \times \mathbb{R}^d \setminus \{0\})$.

The map $\partial_z \sigma(z)(x, \xi)$ is a classical symbol on $U \times \mathbb{R}^d$ and one obtains:

$$\partial_z \sigma(z)(x, \xi) \simeq \sum_{j \geq 0} \partial_z \sigma_{m(z)-j}(z)(x, \xi).$$

Definition 1.17. *The map $P : \Omega \subset \mathbb{C} \rightarrow \Psi DO(U)$ is said to be holomorphic if it has the decomposition*

$$P(z) = \sigma(z)(\cdot, D) + R(z)$$

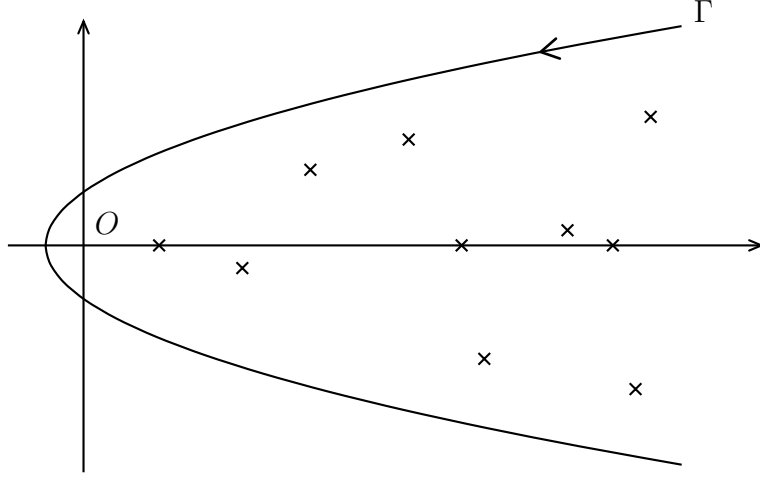
(see definition (1)) where $\sigma : \Omega \rightarrow S(U \times \mathbb{R}^d)$ and $R : \Omega \rightarrow C^\infty(U \times U)$ are holomorphic.

As a consequence, there exists a holomorphic map from Ω into $\Psi DO(M, E)$ with a holomorphic product (when M is compact).

Example 1.18. *Elliptic operators:*

Recall that $P \in \Psi DO^m(U)$, $m \in \mathbb{C}$, is elliptic if there exist strictly positive continuous functions c and C on U such that $|\sigma^P(x, \xi)| \geq c(x) |\xi|^m$ for $|\xi| \geq C(x)$, $x \in U$. This essentially means that P is invertible modulo smoothing operators. More generally, $P \in \Psi DO^m(M, E)$ is elliptic if its local expression in each coordinate chart is elliptic.

Let $Q \in \Psi DO^m(M, E)$ with $\Re(m) > 0$. We assume that M is compact and Q is elliptic. Thus Q has a discrete spectrum and we suppose $\text{Spectrum}(Q) \cap \mathbb{R}^- = \emptyset$. Since we want to integrate in \mathbb{C} , we assume that there exists a curve Γ coming from $+\infty$ along the real axis in the upper half plane, turns around the origin and goes back to infinity in the lower half plane whose interior contains the spectrum of Q . The curve Γ must avoid branch points of λ^z at $z = 0$.



When $\Re(s) < 0$, $Q^s := \frac{1}{i2\pi} \int_{\Gamma} \lambda^s (\lambda - Q)^{-1} d\lambda$ makes sense as operator on $L^2(M, E)$. Actually, $Q^s \in \Psi DO^{ms}(M, E)$ and $(\lambda - Q)^{-1} = \sigma(\lambda)(\cdot, D) + R(\lambda)$ where $R(\lambda)$ is a regularizing operator and $\sigma(\lambda)(\cdot, D)$ has a symbol smooth in λ such that $\sigma(\lambda)(x, \xi) \simeq \sum_{j \geq 0} a_{-m-j}(\lambda, x, \xi)$ with $a_n(\lambda, x, \xi)$ homogeneous of degree n in $(\lambda^{1/m}, \xi)$.

The map $s \rightarrow Q^s$ is a one-parameter group containing $Q^0 = 1$ and $Q^1 = Q$ which is holomorphic on $\Re(s) \leq 0$.

We want to integrate symbols, so we will need the set S_{int} of integrable symbols. Using same type of arguments as in Proposition 1.12 and Lemma 1.15, one proves

Proposition 1.19. *Let*

$$L : \sigma \in S_{int}^{\mathbb{Z}}(\mathbb{R}^d) \rightarrow L(\sigma) := \check{\sigma}(0) = \frac{1}{(2\pi)^d} \int_{\mathbb{R}^d} \sigma(\xi) d\xi.$$

Then L has a unique holomorphic extension \tilde{L} on $S^{\mathbb{C} \setminus \mathbb{Z}}(\mathbb{R}^d)$.

Moreover, when $\sigma(\xi) \simeq \sum_j \sigma_{m-j}(\xi)$, $m \in \mathbb{C} \setminus \mathbb{Z}$,

$$\tilde{L}(\sigma) = \left(\sigma - \sum_{j \leq N} \tau_{m-j} \right)^{\vee}(0) = \frac{1}{(2\pi)^d} \int_{\mathbb{R}^d} \left(\sigma - \sum_{j \leq N} \tau_{m-j} \right)(\xi) d\xi$$

where m is the order of σ , N is an integer with $N > \Re(m) + d$ and τ_{m-j} is the extension of σ_{m-j} of Lemma 1.15.

\tilde{L} is holomorphic extension of L on $S^{\mathbb{C}\setminus\mathbb{Z}}(\mathbb{R}^d)$ which is unique since every element of $S^{\mathbb{C}\setminus\mathbb{Z}}(\mathbb{R}^d)$ is arcwise connected to $S_{int}(\mathbb{R}^d)$ via a holomorphic path within $S^{\mathbb{C}\setminus\mathbb{Z}}(\mathbb{R}^d)$.

This result has an important consequence here:

Corollary 1.20. *If $\sigma : \mathbb{C} \rightarrow S(\mathbb{R}^d)$ is holomorphic and $\text{order}(\sigma(s)) = s$, then $\tilde{L}(\sigma(s))$ is meromorphic with at most simple poles on \mathbb{Z} and for $p \in \mathbb{Z}$,*

$$\text{Res}_{s=p} \tilde{L}(\sigma(s)) = -\frac{1}{(2\pi)^d} \int_{\mathbb{S}^{d-1}} \sigma_{-d}(p)(\xi) d\xi.$$

Proof. Using Lemma 1.15, one proves that if $m(s)$ is holomorphic near $m(s) = p$, then $\tilde{L}(\sigma(s))$ is meromorphic near p .

Now we look at the singularity near $p \in \mathbb{Z}$. In the half plane $\{\Re(s) < p\}$, only the infrared part of $\tau_{m-j}(s)$ is a problem since its ultraviolet part is holomorphic. For $0 \leq j \leq p + m$ and $\Re(s) < p$, $\sigma_{s-j}(s)(\xi)$ is integrable near 0 thus defines its unique extension $\tau_{s-j}(s)$. So, the only possible singularity near $s = p$ could come from

$$\begin{aligned} -\frac{1}{(2\pi)^d} \int_{|\xi| \leq 1} \sigma_{s-j}(s)(\xi) d\xi &= -\frac{1}{(2\pi)^d} \int_0^1 t^{s-j+d-1} dt \int_{|\xi| \leq 1} \sigma_{s-j}(s)\left(\frac{\xi}{|\xi|}\right) d\left(\frac{\xi}{|\xi|}\right) \\ &= -\frac{1}{(2\pi)^d} \frac{1}{s-j+d} \int_{\mathbb{S}^{d-1}} \sigma_{s-j}(s)(\xi) d\xi. \end{aligned}$$

where we used for the first equality $\sigma_{s-j}(s)(\xi) = |\xi|^{s-j} \sigma_{s-j}(s)\left(\frac{\xi}{|\xi|}\right)$. Thus, $\tilde{L}(\sigma(s))$ has at most only simple pole at $s = -d + j$. \square

We are now ready to get the main result of this section which is due to Wodzicki [111, 112].

Definition 1.21. *Let $\mathcal{D} \in \Psi DO(M, E)$ be an elliptic pseudodifferential operator of order 1 on a boundary-less compact manifold M endowed with a vector bundle E .*

Let $\Psi DO_{int}(M, E) := \{Q \in \Psi DO^{\mathbb{C}}(M, E) \mid \Re(\text{order}(Q)) < -d\}$ be the class of pseudodifferential operators whose symbols are in S_{int} , i.e. integrable in the ξ -variable.

In particular, if $P \in \Psi DO_{int}(M, E)$, then its kernel $k^P(x, x)$ is a smooth density on the diagonal of $M \times M$ with values in $\text{End}(E)$.

For $P \in \Psi DO^{\mathbb{Z}}(M, E)$, define

$$WRes P := \text{Res}_{s=0} \text{Tr} \left(P|\mathcal{D}|^{-s} \right). \quad (11)$$

This makes sense because:

Theorem 1.22. *(i) The map $P \in \Psi DO_{int}(M, E) \rightarrow \text{Tr}(P) \in \mathbb{C}$ has a unique analytic extension on $\Psi DO^{\mathbb{C}\setminus\mathbb{Z}}(M, E)$.*

(ii) If $P \in \Psi DO^{\mathbb{Z}}(M, E)$, the map: $s \in \mathbb{C} \rightarrow \text{Tr} \left(P|\mathcal{D}|^{-s} \right)$ has at most simple poles on \mathbb{Z} and

$$WRes P = - \int_M c_P(x) |dx| \quad (12)$$

is independent of \mathcal{D} . Recall (see Theorem 1.14) that $c_P(x) = \frac{1}{(2\pi)^d} \int_{\mathbb{S}^{d-1}} \text{tr} \left(\sigma_{-d}^P(x, \xi) \right) d\xi$.

(iii) $WRes$ is a trace on the algebra $\Psi DO^{\mathbb{Z}}(M, E)$.

Proof. (i) The map $s \rightarrow \text{Tr} \left(P|D|^{-s} \right)$ is holomorphic on \mathbb{C} and connect $P \in \Psi DO^{\mathbb{C}\mathbb{Z}}(M, E)$ to the set $\Psi DO_{int}(M, E)$ within $\Psi DO^{\mathbb{C}\mathbb{Z}}(M, E)$, so a analytic extension is necessarily unique.

(ii) one apply the above machinery:

(1) Notice that Tr is holomorphic on smoothing operator, so, using a partition of unity, we can reduce to a local study of scalar ΨDO 's.

(2) First, fix $s = 0$. We are interested in the function $L_\phi(\sigma) := \text{Tr} \left(\phi \sigma(x, D) \right)$ with $\sigma \in S_{int}(U \times \mathbb{R}^d)$ and $\phi \in C^\infty(U)$. For instance, if $P = \sigma(\cdot, D)$,

$$\text{Tr}(\phi P) = \int_U \phi(x) k^P(x, x) |dx| = \frac{1}{(2\pi)^d} \int_U \phi(x) \sigma(x, \xi) d\xi |dx| = \int_U \phi(x) L(\sigma(x, \cdot)) |dx|,$$

so one extends L_ϕ to $S^{\mathbb{C}\mathbb{Z}}(U \times \mathbb{R}^d)$ with Proposition 1.19 via $\tilde{L}_\phi(\sigma) = \int_U \phi(x) \tilde{L}_\phi(\sigma(x, \cdot)) |dx|$.

(3) If now $\sigma(x, \xi) = \sigma(s)(x, \xi)$ depends holomorphically on s , we get uniform bounds in x , thus we get, via Lemma 1.15 applied to $\tilde{L}_\phi(\sigma(s)(x, \cdot))$ uniformly in x , yielding a natural extension to $\tilde{L}_\phi(\sigma(s))$ which is holomorphic on $\mathbb{C} \setminus \mathbb{Z}$.

When $\text{order}(\sigma(s)) = s$, the map $\tilde{L}_\phi(\sigma(s))$ has at most simple poles on \mathbb{Z} and for each $p \in \mathbb{Z}$, $\text{Res}_{s=p} \tilde{L}_\phi(\sigma(s)) = -\frac{1}{(2\pi)^d} \int_U \int_{\mathbb{S}^{d-1}} \phi(x) \sigma_{-d}(p)(x, \xi) d\xi |dx| = -\int_U \phi(x) c_{P_p}(x) |dx|$ where we used (9) with $P = Op(\sigma_p(x, \xi))$.

(4) In the general case, we get a unique meromorphic extension of the usual trace Tr on $\Psi DO^{\mathbb{Z}}(M, E)$ that we still denoted by Tr .

When $P : \mathbb{C} \rightarrow \Psi DO^{\mathbb{Z}}(M, E)$ is meromorphic with $\text{order}(P(s)) = s$, then $\text{Tr}(P(s))$ has at most poles on \mathbb{Z} and $\text{Res}_{s=p} \text{Tr}(P(s)) = -\int_M c_{P(p)}(x) |dx|$ for $p \in \mathbb{Z}$. So we get the claim for the family

$$P(s) := P|D|^{-s}.$$

(iii) Let $P_1, P_2 \in \Psi DO^{\mathbb{Z}}(M, E)$. Since Tr is a trace on $\Psi DO^{\mathbb{C}\mathbb{Z}}(M, E)$, we get by (i), $\text{Tr}(P_1 P_2 |D|^{-s}) = \text{Tr}(P_2 |D|^{-s} P_1)$. Moreover

$$WRes(P_1 P_2) = \text{Res}_{s=0} \text{Tr}(P_2 |D|^{-s} P_1) = \text{Res}_{s=0} \text{Tr}(P_2 P_1 |D|^{-s}) = WRes(P_2 P_1)$$

where for the second equality we used (12) so the residue depends only of the value of $P(s)$ at $s = 0$. \square

Note that $WRes$ is invariant by diffeomorphism:

$$\text{if } \phi \in \text{Diff}(M), WRes(P) = WRes(\phi_* P) \quad (13)$$

which follows from (10). The next result is due to Guillemin and Wodzicki.

Corollary 1.23. *The Wodzicki residue $WRes$ is the only trace (up to multiplication by a constant) on the algebra $\Psi DO^{-m}(M, E)$, $m \in \mathbb{N}$, when M is connected and $d \geq 2$.*

Proof. The restriction to $d \geq 2$ is used only in the part 3) below. When $d = 1$, T^*M is disconnected and they are two residues.

1) On symbols, derivatives are commutators:

$$[x^j, \sigma] = i\partial_{\xi_j} \sigma, \quad [\xi_j, \sigma] = -i\partial_{x^j} \sigma.$$

2) If $\sigma_{-d}^P = 0$, then $\sigma^P(x, \xi)$ is a finite sum of commutators of symbols:
When $\sigma^P \simeq \sum_j \sigma_{m-j}^P$ with $m = \text{order}(P)$, by Euler's theorem,

$$\sum_{k=1}^d \xi_k \partial_{\xi_k} \sigma_{m-j}^P = (m-j) \sigma_{m-j}^P$$

(this is false for $m = j!$) and

$$\sum_{k=1}^d [x^k, \xi_k \sigma_{m-j}^P] = i \sum_{k=1}^d \partial_{\xi_k} \xi_k \sigma_{m-j}^P = i(m-j+d) \sigma_{m-j}^P.$$

So $\sigma^P = \sum_{k=1}^d [\xi_k \tau, x^k]$ (where $\tau \simeq i \sum_{j \geq 0} \frac{1}{m-j+d} \sigma_{m-j}^P$ and here we need for $m-j = -d$ that $\sigma_{-d}^P = 0!$).

Let T be another trace on $\Psi DO^{\mathbb{Z}}(M, E)$. Then $T(P)$ depends only on σ_{-d}^P because $T([\cdot, \cdot]) = 0$.

3) We have $\int_{\mathbb{S}^{d-1}} \sigma_{-d}^P(x, \xi) d|\xi| = 0$ if and only if σ_{-d}^P is sum of derivatives:

The if part is direct (less than more!).

Only if part: σ_{-d}^P is orthogonal to constant functions on the sphere \mathbb{S}^{d-1} and these are kernels of the Laplacian: $\Delta_{\mathbb{S}} f = 0 \iff df = 0 \iff f = \text{cst}$. Thus $\Delta_{\mathbb{S}^{d-1}} h = \sigma_{-d}^P|_{\mathbb{S}^{d-1}}$ has a solution h on \mathbb{S}^{d-1} . If $\tilde{h}(\xi) := |\xi|^{-d+2} h\left(\frac{\xi}{|\xi|}\right)$ is its extension to $\mathbb{R}^d \setminus \{0\}$, then we get $\Delta_{\mathbb{R}^d} \tilde{h}(\xi) = |\xi| \sigma_{-d}^P\left(\frac{\xi}{|\xi|}\right) = \sigma_{-d}^P(\xi)$ because $\Delta_{\mathbb{R}^d} = r^{1-d} \partial_r (r^{d-1} \partial_r) + r^{-2} \Delta_{\mathbb{S}^{d-1}}$. This means that \tilde{h} is a symbol of order $d-2$ and $\partial_{\xi} \tilde{h}$ is a symbol of order $d-1$. As a consequence, $\sigma_{-d}^P = \sum_{k=1}^d \partial_{\xi_k}^2 \tilde{h} = -i \sum_{k=1}^d [\partial_{\xi_k} \tilde{h}, x^k]$ is a sum of commutators.

4) End of proof:

$\sigma_{-d}^P(x, \xi) - \frac{|\xi|^{-d}}{\text{Vol}(\mathbb{S}^{d-1})} c_P(x)$ is a symbol of order $-d$ with zero integral, thus is a sum of commutators by 3) and $T(P) = T(\text{Op}(|\xi|^{-d} c_P(x)))$ for all $T \in \Psi DO^{\mathbb{Z}}(M, E)$. In other words, the map $\mu : f \in C_c^{\infty}(U) \rightarrow T(\text{Op}(f|\xi|^{-d}))$ is linear, continuous and satisfies $\mu(\partial_{x^k} f) = 0$ because $\partial_{x^k}(f) |\xi|^{-d}$ is a commutator if f has a compact support and U is homeomorphic to \mathbb{R}^d . As a consequence, μ is a multiple of the Lebesgue integral:

$$T(P) = \mu(c_P(x)) = c \int_M c_P(x) |dx| = c \text{WRes}(P). \quad \square$$

Example 1.24. *Laplacian on a manifold M : Let M be a compact Riemannian manifold of dimension d and Δ be the scalar Laplacian which is a differential operator of order 2. Then*

$$\text{WRes}\left((1 + \Delta)^{-d/2}\right) = \text{Vol}\left(\mathbb{S}^{d-1}\right) = \frac{2\pi^{d/2}}{\Gamma(d/2)}.$$

Proof. $(1 + \Delta)^{-d/2} \in \Psi DO(M)$ has order $-d$ and its principal symbol $\sigma_{-d}^{(1+\Delta)^{-d/2}}$ satisfies

$$\sigma_{-d}^{(1+\Delta)^{-d/2}}(x, \xi) = -\left(g_x^{ij} \xi_i \xi_j\right)^{-d/2} = -\|\xi\|_x^{-d}.$$

So (12) gives

$$\begin{aligned} \text{WRes}\left((1 + \Delta)^{-d/2}\right) &= \int_M |dx| \int_{\mathbb{S}^{d-1}} \|\xi\|_x^{-d} d\xi = \int_M |dx| \sqrt{\det g_x} \text{Vol}\left(\mathbb{S}^{d-1}\right) \\ &= \text{Vol}\left(\mathbb{S}^{d-1}\right) \int_M |d\text{vol}_g| = \text{Vol}\left(\mathbb{S}^{d-1}\right). \end{aligned} \quad \square$$

2 Dixmier trace

References for this section: [34, 50, 69, 89, 106, 107].

The trace on the operators on a Hilbert space \mathcal{H} has an interesting property, it is *normal*. Recall first that Tr acting on $\mathcal{B}(\mathcal{H})$ is a particular case of a weight ω acting on a von Neumann algebra \mathcal{M} : it is a homogeneous additive map from positive elements $\mathcal{M}^+ := \{aa^* \mid a \in \mathcal{M}\}$ to $[0, \infty]$.

A state is a weight $\omega \in \mathcal{M}^*$ (so $\omega(a) < \infty, \forall a \in \mathcal{M}$) such that $\omega(1) = 1$.

A trace is a weight such that $\omega(aa^*) = \omega(a^*a)$ for all $a \in \mathcal{M}$.

Definition 2.1. *A weight ω is normal if $\omega(\sup_{\alpha} a_{\alpha}) = \sup_{\alpha} \omega(a_{\alpha})$ whenever $(a_{\alpha}) \subset \mathcal{M}^+$ is an increasing bounded net.*

This is equivalent to say that ω is lower semi-continuous with respect to the σ -weak topology.

Lemma 2.2. *The usual trace Tr is normal on $\mathcal{B}(\mathcal{H})$.*

Remark that the net $(a_{\alpha})_{\alpha}$ converges in $\mathcal{B}(\mathcal{H})$ and this property looks innocent since a trace preserves positivity.

Nevertheless it is natural to address the question: are all traces (in particular on an arbitrary von Neumann algebra) normal? In 1966, Dixmier answered by the negative [35] by exhibiting non-normal, say singular, traces. Actually, his motivation was to answer the following related question: is any trace ω on $\mathcal{B}(\mathcal{H})$ proportional to the usual trace on the set where ω is finite?

The aim of this section is first to define this Dixmier trace, which essentially means $\text{Tr}_{Dix}(T) \text{ " = " } \lim_{N \rightarrow \infty} \frac{1}{\log N} \sum_{n=0}^N \mu_n(T)$, where the $\mu_n(T)$ are the singular values of T ordered in decreasing order and then to relate this to the Wodzicki trace. It is a non-normal trace on some set that we have to identify. Naturally, the reader can feel the link with the Wodzicki trace via Proposition 1.13. We will see that if $P \in \Psi DO^{-d}(M)$ where M is a compact Riemannian manifold of dimension d , then,

$$\text{Tr}_{Dix}(P) = \frac{1}{d} WRes(P) = \frac{1}{d} \int_M \int_{S^*M} \sigma_{-d}^P(x, \xi) d\xi |dx|$$

where S^*M is the cosphere bundle on M .

The physical motivation is quite essential: We know how $\sum_{n \in \mathbb{N}^*} \frac{1}{n}$ diverges and this is related to the fact the electromagnetic or Newton gravitational potentials are in $\frac{1}{r}$ which has the same singularity (in one-dimension as previous series). Actually, this (logarithmic-type) divergence appears everywhere in physics and explains the widely use of the Riemann zeta function $\zeta : s \in \mathbb{C} \rightarrow \sum_{n \in \mathbb{N}^*} \frac{1}{n^s}$. This is also why we have already seen a logarithmic obstruction in Theorem 1.14 and define a zeta function associated to a pseudodifferential operator P by $\zeta_P(s) = \text{Tr}(P|\mathcal{D}|^{-s})$ in (11).

We now have a quick review on the main properties of singular values of an operator.

2.1 Singular values of compact operators

In noncommutative geometry, infinitesimals correspond to compact operators: for $T \in \mathcal{K}(\mathcal{H})$ (compact operators), define for $n \in \mathbb{N}$

$$\mu_n(T) := \inf \{ \|T|_{E^\perp}\| \mid E \text{ subspace of } \mathcal{H} \text{ with } \dim(E) = n \}.$$

This could look strange but actually, by mini-max principle, $\mu_n(T)$ is nothing else than the $(n+1)$ th of eigenvalues of $|T|$ sorted in decreasing order. Since $\lim_{n \rightarrow \infty} \mu_n(T) = 0$, for any $\epsilon > 0$, there exists a finite-dimensional subspace E_ϵ such that $\|T|_{E_\epsilon^\perp}\| < \epsilon$ and this property being equivalent to T compact, T deserves the name of infinitesimal.

Moreover, we have following properties:

$$\begin{aligned} \mu_n(T) &= \mu_n(T^*) = \mu_n(|T|). \\ T \in \mathcal{L}^1(\mathcal{H}) \text{ (meaning } \|T\|_1 := \text{Tr}(|T|) < \infty) &\iff \sum_{n \in \mathbb{N}} \mu_n(T) < \infty. \\ \mu_n(A T B) &\leq \|A\| \mu_n(T) \|B\| \text{ when } A, B \in \mathcal{B}(\mathcal{H}). \\ \mu_N(U T U^*) &= \mu_N(T) \text{ when } U \text{ is a unitary.} \end{aligned}$$

Definition 2.3. For $T \in \mathcal{K}(\mathcal{H})$, the partial trace of order $N \in \mathbb{N}$ is $\sigma_N(T) := \sum_{n=0}^N \mu_n(T)$.

Remark that $\|T\| \leq \sigma_N(T) \leq N\|T\|$ which implies $\sigma_n \simeq \|\cdot\|$ on $\mathcal{K}(\mathcal{H})$. Then

$$\begin{aligned} \sigma_N(T_1 + T_2) &\leq \sigma_N(T_1) + \sigma_N(T_2), \\ \sigma_{N_1}(T_1) + \sigma_{N_2}(T_2) &\leq \sigma_{N_1+N_2}(T_1 + T_2) \text{ when } T_1, T_2 \geq 0. \end{aligned} \quad (14)$$

The proof of the sub-additivity is based on the fact that σ_N is a norm on $\mathcal{K}(\mathcal{H})$. Moreover

$$T \geq 0 \implies \sigma_N(T) = \sup\{ \text{Tr}(T E) \mid E \text{ subspace of } \mathcal{H} \text{ with } \dim(E) = n \}.$$

which implies $\sigma_N(T) = \sup\{ \text{Tr}(\|T E\|_1 \mid \dim(E) = n) \}$ and gives the second inequality.

The norm σ_N can be decomposed:

$$\sigma_N(T) = \inf\{ \|x\|_1 + N\|y\| \mid T = x + y \text{ with } x \in \mathcal{L}^1(\mathcal{H}), y \in \mathcal{K}(\mathcal{H}) \}.$$

In fact if $\tilde{\sigma}_N$ is the right hand-side, then the sub-additivity gives $\tilde{\sigma}_N \geq \sigma_N(T)$. To get equality, let $\xi_n \in \mathcal{H}$ be such that $|T|\xi_n = \mu_n(T)\xi_n$ and define $x_N := (|T| - \mu_N(T))E_N$, $y_N := \mu_N(T)E_N + |T|(1 - E_N)$ where $E_N := \sum_{n < N} |\xi_n\rangle\langle\xi_n|$. If $T = U|T|$ is the polar decomposition of T , then $T = Ux_N + Uy_N$ is a claimed decomposition of T and

$$\tilde{\sigma}_N(T) \leq \|Ux_N\|_1 + N\|Uy_N\| \leq \|x_N\|_1 + N\|y_N\| \leq \sum_{n < N} (\mu_n(T) - N\mu_n(T)) + N\mu_N(T) \leq \sigma_N(T).$$

This justifies a continuous approach with the

Definition 2.4. The partial trace of T of order $\lambda \in \mathbb{R}^+$ is

$$\sigma_\lambda(T) := \inf\{ \|x\|_1 + \lambda\|y\| \mid T = x + y \text{ with } x \in \mathcal{L}^1(\mathcal{H}), y \in \mathcal{K}(\mathcal{H}) \}.$$

It interpolates between two consecutive integers since the map: $\lambda \rightarrow \sigma_\lambda(T)$ is concave for $T \in \mathcal{K}(\mathcal{H})$ and moreover, it is affine between N and $N+1$ because

$$\sigma_\lambda(T) = \sigma_N(T) + (\lambda - N)\sigma_N(T), \text{ where } N = [\lambda]. \quad (15)$$

Thus, as before,

$$\sigma_{\lambda_1}(T_1) + \sigma_{\lambda_2}(T_2) = \sigma_{\lambda_1 + \lambda_2}(T_1 + T_2), \text{ for } \lambda_1, \lambda_2 \in \mathbb{R}^+, 0 \leq T_1, T_2 \in \mathcal{K}(\mathcal{H}).$$

We define a real interpolate space between $\mathcal{L}^1(\mathcal{H})$ and $\mathcal{K}(\mathcal{H})$ by

$$\mathcal{L}^{1,\infty} := \{ T \in \mathcal{K}(\mathcal{H}) \mid \|T\|_{1,\infty} := \sup_{\lambda \geq e} \frac{\sigma_\lambda(T)}{\log \lambda} < \infty \}.$$

If $\mathcal{L}^p(\mathcal{H})$ is the ideal of operators T such that $\text{Tr}(|T|^p) < \infty$, so $\sigma_\lambda(T) = \mathcal{O}(\lambda^{1-1/p})$, we have naturally

$$\begin{aligned} \mathcal{L}^1(\mathcal{H}) &\subset \mathcal{L}^{1,\infty} \subset \mathcal{L}^p(\mathcal{H}) \text{ for } p > 1, \\ \|T\| &\leq \|T\|_{1,\infty} \leq \|T\|_1. \end{aligned} \tag{16}$$

Lemma 2.5. $\mathcal{L}^{1,\infty}$ is a C^* -ideal of $\mathcal{B}(\mathcal{H})$ for the norm $\|\cdot\|_{1,\infty}$. Moreover, it is equal to the Macaev ideal

$$\mathcal{L}^{1,+} := \{ T \in \mathcal{K}(\mathcal{H}) \mid \|T\|_{1,+} := \sup_{N \geq 2} \frac{\sigma_N(T)}{\log(N)} < \infty \}.$$

Proof. $\|\cdot\|_{1,\infty}$ is a norm as supremum of norms. By (15),

$$\sup_{\rho \geq e} \frac{\sigma_\rho(T)}{\log \rho} \leq \sup_{N \geq 2} \sup_{0 \leq \alpha \leq 1} \frac{\sum_{n=0}^{N-1} \mu_N(T) + \alpha \mu_N(T)}{\log(N + \alpha)}$$

and $\mathcal{L}^{1,+}$ is a left and right ideal of $\mathcal{B}(\mathcal{H})$ since $\|ATB\|_{1,\infty} \leq \|A\| \|T\|_{1,\infty} \|B\|$ for every $A, B \in \mathcal{B}(\mathcal{H})$, $T \in \mathcal{L}^{1,\infty}$, and moreover $\|T\|_{1,\infty} = \|T^*\|_{1,\infty} = \||T|\|_{1,\infty}$.

This ideal $\mathcal{L}^{1,\infty}$ is closed for $\|\cdot\|_{1,\infty}$: this follows from a $3-\epsilon$ argument since Cauchy sequences for $\|\cdot\|_{1,\infty}$ are Cauchy sequences for each norm σ_λ which are equivalent to $\|\cdot\|$. \square

Despite this result, the reader should notice that $\|\cdot\|_{1,\infty} \neq \|\cdot\|_{1,+}$ since the norms are only equivalent.

2.2 Dixmier trace

We begin with a Cesàro mean of $\frac{\sigma_\rho(T)}{\log \rho}$ with respect of the Haar measure of the group \mathbb{R}_+^* :

Definition 2.6. For $\lambda \geq e$ and $T \in \mathcal{K}(\mathcal{H})$, let

$$\tau_\lambda(T) := \frac{1}{\log \lambda} \int_e^\lambda \frac{\sigma_\rho(T)}{\log \rho} \frac{d\rho}{\rho}.$$

Clearly, $\sigma_\rho(T) \leq \log \rho \|T\|_{1,\infty}$ and $\tau_\lambda(T) \leq \|T\|_{1,\infty}$, thus the map: $\lambda \rightarrow \tau_\lambda(T)$ is in $C_b([e, \infty])$. It is not additive on $\mathcal{L}^{1,\infty}$ but this defect is under control:

$$\tau_\lambda(T_1 + T_2) - \tau_\lambda(T_1) - \tau_\lambda(T_2) \underset{\lambda \rightarrow \infty}{\simeq} \mathcal{O}\left(\frac{\log(\log \lambda)}{\log \lambda}\right), \text{ when } 0 \leq T_1, T_2 \in \mathcal{L}^{1,\infty}.$$

More precisely, using previous results, one get

Lemma 2.7.

$$|\tau_\lambda(T_1 + T_2) - \tau_\lambda(T_1) - \tau_\lambda(T_2)| \leq \left(\frac{\log 2(2 + \log \log \lambda)}{\log \lambda}\right) \|T_1 + T_2\|_{1,\infty}, \text{ when } T_1, T_2 \in \mathcal{L}_+^{1,\infty}.$$

Proof. By the sub-additivity of σ_ρ , $\tau_\lambda(T_1 + T_2) \leq \tau_\lambda(T_1) + \tau_\lambda(T_2)$ and thanks to (14), we get $\sigma_\rho(T_1) + \sigma_\rho(T_2) \leq \sigma_{2\rho}(T_1 + T_2)$. Thus

$$\tau_\lambda(T_1) + \tau_\lambda(T_2) \leq \frac{1}{\log \lambda} \int_e^\lambda \frac{\sigma_{2\rho}(T_1 + T_2)}{\log \rho} \frac{d\rho}{\rho} \leq \frac{1}{\log \lambda} \int_{2e}^{2\lambda} \frac{\sigma_\rho(T_1 + T_2)}{\log \rho/2} \frac{d\rho}{\rho}$$

Hence, $(\log \lambda) |\tau_\lambda(T_1 + T_2) - \tau_\lambda(T_1) - \tau_\lambda(T_2)| \leq \epsilon + \epsilon'$ with

$$\begin{aligned}\epsilon &:= \int_e^\lambda \frac{\sigma_\rho(T_1+T_2)}{\log \rho} \frac{d\rho}{\rho} - \int_{2e}^{2\lambda} \frac{\sigma_\rho(T_1+T_2)}{\log \rho/2} \frac{d\rho}{\rho}, \\ \epsilon' &:= \int_{2e}^{2\lambda} \sigma_\rho(T_1 + T_2) \left(\frac{1}{\log \rho/2} - \frac{1}{\log \rho} \right) \frac{d\rho}{\rho}.\end{aligned}$$

By triangular inequality and the fact that $\sigma_\rho(T_1 + T_2) \leq \log \rho \|T_1 + T_2\|_{1,\infty}$ when $\rho \geq e$,

$$\begin{aligned}\epsilon &\leq \int_e^{2e} \frac{\sigma_\rho(T_1+T_2)}{\log \rho} \frac{d\rho}{\rho} + \int_\lambda^{2\lambda} \frac{\sigma_\rho(T_1+T_2)}{\log \rho} \frac{d\rho}{\rho} \\ &\leq \|T_1 + T_2\|_{1,\infty} \left(\int_e^{2e} \frac{d\rho}{\rho} + \int_\lambda^{2\lambda} \frac{d\rho}{\rho} \right) \leq 2 \log(2) \|T_1 + T_2\|_{1,\infty}.\end{aligned}$$

Moreover,

$$\begin{aligned}\epsilon' &\leq \|T_1 + T_2\|_{1,\infty} \int_{2e}^{2\lambda} \log \rho \left(\frac{1}{\log \rho/2} - \frac{1}{\log \rho} \right) \frac{d\rho}{\rho} \leq \|T_1 + T_2\|_{1,\infty} \int_{2e}^{2\lambda} \frac{\log 2}{\log \rho/2} \frac{d\rho}{\rho} \\ &\leq \|T_1 + T_2\|_{1,\infty} \log(2) \log(\log \lambda).\end{aligned} \quad \square$$

The Dixmier's idea was to force additivity: since the map $\lambda \rightarrow \tau_\lambda(T)$ is in $C_b([e, \infty])$ and $\lambda \rightarrow \left(\frac{\log 2(2+\log \log \lambda)}{\log \lambda} \right)$ is in $C_0([e, \infty])$, let us consider the C^* -algebra

$$\mathcal{A} := C_b([e, \infty]) / C_0([e, \infty]).$$

If $[\tau(T)]$ is the class of the map $\lambda \rightarrow \tau_\lambda(T)$ in \mathcal{A} , previous lemma shows that $[\tau] : T \rightarrow [\tau(T)]$ is additive and positive homogeneous from $\mathcal{L}_+^{1,\infty}$ into \mathcal{A} satisfying $[\tau(UTU^*)] = [\tau(T)]$ for any unitary U .

Now let ω be a state on \mathcal{A} , namely a positive linear form on \mathcal{A} with $\omega(1) = 1$.

Then, $\omega \circ [\tau(\cdot)]$ is a tracial weight on $\mathcal{L}_+^{1,\infty}$ (a map from $\mathcal{L}_+^{1,\infty}$ to \mathbb{R}^+ which is additive, homogeneous and invariant under $T \rightarrow UTU^*$). Since $\mathcal{L}^{1,\infty}$ is a C^* -ideal of $\mathcal{B}(\mathcal{H})$, each of its element is generated by (at most) four positive elements, and this map can be extended to a map $\omega \circ [\tau(\cdot)] : T \in \mathcal{L}^{1,\infty} \rightarrow \omega([\tau(T)]) \in \mathbb{C}$ such that $\omega([\tau(T_1 T_2)]) = \omega([\tau(T_2 T_1)])$ for $T_1, T_2 \in \mathcal{L}^{1,\infty}$. This leads to the following

Definition 2.8. *The Dixmier trace Tr_ω associated to a state ω on $\mathcal{A} := C_b([e, \infty]) / C_0([e, \infty])$ is*

$$\text{Tr}_\omega(\cdot) := \omega \circ [\tau(\cdot)].$$

Theorem 2.9. *Tr_ω is a trace on $\mathcal{L}^{1,\infty}$ which depends only on the locally convex topology of \mathcal{H} , not of its scalar product.*

Proof. We already know that Tr_ω is a trace.

If $\langle \cdot, \cdot \rangle'$ is another scalar product on \mathcal{H} giving the same topology as $\langle \cdot, \cdot \rangle$, then there exist an invertible $U \in \mathcal{B}(\mathcal{H})$ with $\langle \cdot, \cdot \rangle' = \langle U \cdot, U \cdot \rangle$. Let \mathcal{H}' be the Hilbert space for $\langle \cdot, \cdot \rangle'$ and Tr'_ω be the associated Dixmier trace to a given state ω . Since the singular value of $U^{-1}TU \in \mathcal{K}_+(\mathcal{H}')$ are the same of $T \in \mathcal{K}_+(\mathcal{H})$, we get $\mathcal{L}^{1,\infty}(\mathcal{H}') = \mathcal{L}^{1,\infty}(\mathcal{H})$ and

$$\text{Tr}'_\omega(T) = \text{Tr}_\omega(U^{-1}TU) = \text{Tr}_\omega(T) \text{ for } T \in \mathcal{L}_+^{1,\infty}. \quad \square$$

Two important points:

1) Note that $\text{Tr}_\omega(T) = 0$ if $T \in \mathcal{L}^1(\mathcal{H})$ and more generally all Dixmier traces vanish on the closure for the norm $\|\cdot\|_{1,\infty}$ of the ideal of finite rank operators. In particular, Dixmier traces are not normal.

2) The C^* -algebra \mathcal{A} is not separable, so it is impossible to exhibit any state ω ! Despite the inclusions (16) and the fact that the $\mathcal{L}^p(\mathcal{H})$ are separable ideals for $p \geq 1$, $\mathcal{L}^{1,\infty}$ is not a separable.

Moreover, as for Lebesgue integral, there are sets which are not measurable. For instance, a function $f \in C_b([e, \infty))$ has a limit $\ell = \lim_{\lambda \rightarrow \infty} f(\lambda)$ if and only if $\ell = \omega(f)$ for all state ω .

Definition 2.10. *The operator $T \in \mathcal{L}^{1,\infty}$ is said to be measurable if $\text{Tr}_\omega(T)$ is independent of ω . In this case, Tr_ω is denoted Tr_{Dix} .*

Lemma 2.11. *The operator $T \in \mathcal{L}^{1,\infty}$ is measurable and $\text{Tr}_\omega(T) = \ell$ if and only if the map $\lambda \in \mathbb{R}^+ \rightarrow \tau_\lambda(T) \in \mathcal{A}$ converges at infinity to ℓ .*

Proof. If $\lim_{\tau \rightarrow \infty} \tau_\lambda(T) = \ell$, then $\text{Tr}_\omega(T) = \omega(\tau(T)) = \omega(\ell) = \ell \omega(1) = \ell$.

Conversely, assume T is measurable and $\ell = \text{Tr}_\omega(T)$ for any state ω . Then we get, $\omega(\tau(T) - \ell) = \text{Tr}_\omega(T) - \ell = 0$. Since the set of states separate the points of \mathcal{A} , $\tau(T) = \ell$ and $\lim_{\tau \rightarrow \infty} \tau_\lambda(T) = \ell$. \square

After Dixmier, the singular (i.e. non normal) traces have been deeply investigated, see for instance the recent [73, 75, 76], but we do not enter into this framework and technically, we just make the following characterization of measurability:

Remark 2.12. *If $T \in \mathcal{K}_+(\mathcal{H})$, then T is measurable if and only if $\lim_{N \rightarrow \infty} \frac{1}{\log N} \sum_{n=0}^N \mu_n(T)$ exists.*

Actually, if $\ell = \lim_{N \rightarrow \infty} \frac{1}{\log N} \sum_{n=0}^N \mu_n(T)$ since $\text{Tr}_\omega(T) = \ell$ for any ω , so $\text{Tr}_{Dix} = \ell$ and the converse is proved in [74].

Example 2.13. *Computation of the Dixmier trace of the inverse Laplacian on the torus:*

Let $\mathbb{T}^d = \mathbb{R}^d / 2\pi\mathbb{Z}^d$ be the d -dimensional torus and $\Delta = -\sum_{i=1}^d \partial_{x_i}^2$ be the scalar Laplacian seen as unbounded operator on $\mathcal{H} = L^2(\mathbb{T}^d)$. We want to compute $\text{Tr}_\omega((1 + \Delta)^{-p})$ for $\frac{d}{2} \leq p \in \mathbb{N}^*$. We use $1 + \Delta$ to avoid the kernel problem with the inverse. As the following proof shows, 1 can be replaced by any $\epsilon > 0$ and the result does not depends on ϵ .

Notice that the functions $e_k(x) := \frac{1}{2\pi} e^{ik \cdot x}$ with $x \in \mathbb{T}^d$, $k \in \mathbb{Z}^d = (\mathbb{T}^d)^*$ form a basis of \mathcal{H} of eigenvectors: $\Delta e_k = |k|^2 e_k$. Moreover, for $t \in \mathbb{R}_+$,

$$e^t \text{Tr} \left(e^{-t(1+\Delta)} \right) = \sum_{k \in \mathbb{Z}^d} e^{-t|k|^2} = \left(\sum_{k \in \mathbb{Z}} e^{-tk^2} \right)^d.$$

We know that $|\int_{-\infty}^{\infty} e^{-tx^2} dx - \sum_{k \in \mathbb{Z}} e^{-tk^2}| \leq 1$, and since the first integral is $\sqrt{\frac{\pi}{t}}$, we get $e^t \text{Tr} \left(e^{-t(1+\Delta)} \right) \underset{t \downarrow 0^+}{\simeq} \left(\frac{\pi}{t} \right)^{d/2} =: \alpha t^{-d/2}$.

We will use a Tauberian theorem: $\mu_n \left((1 + \Delta)^{-d/2} \right) \underset{n \rightarrow \infty}{\simeq} \left(\alpha \frac{1}{\Gamma(d/2+1)} \right) \frac{1}{n}$, see [55] (one needs to estimates the cardinality of the set $\{k \in \mathbb{Z}^d \mid |k|^2 \leq n\}$, see [50]). Thus

$$\lim_{N \rightarrow \infty} \frac{1}{\log N} \sum_{n=0}^N \mu_n \left((1 + \Delta)^{-d/2} \right) = \frac{\alpha}{\Gamma(d/2+1)} = \frac{\pi^{d/2}}{\Gamma(d/2+1)}.$$

Thus $(1 + \Delta)^{-d/2}$ is measurable and

$$\mathrm{Tr}_{Dix} \left((1 + \Delta)^{-d/2} \right) = \mathrm{Tr}_\omega \left((1 + \Delta)^{-d/2} \right) = \frac{\pi^{d/2}}{\Gamma(d/2+1)}.$$

Since $(1 + \Delta)^{-p}$ is traceable for $p > \frac{d}{2}$, $\mathrm{Tr}_{Dix} \left((1 + \Delta)^{-p} \right) = 0$.

This result has been generalized in Connes' trace theorem [24]:

Theorem 2.14. *Let M be a compact Riemannian manifold of dimension d , E a vector bundle over M and $P \in \Psi DO^{-d}(M, E)$. Then, $P \in \mathcal{L}^{1,\infty}$, is measurable and*

$$\mathrm{Tr}_{Dix}(P) = \frac{1}{d} WRes(P).$$

Proof. Since $WRes$ and Tr_{Dix} are traces on $\Psi DO^{-m}(M, E)$, $m \in \mathbb{N}$, $\mathrm{Tr}_{Dix} = c WRes$ for some constant c using Corollary 1.23. Above example, when compare with Example 1.24, shows that the $c = \frac{1}{d}$. \square

3 Dirac operator

There are several ways to define a Dirac-like operator. The best one is to define Clifford algebras, their representations, the notion of Clifford modules, spin^c structures on orientable manifolds M defined by Morita equivalence between the C^* -algebras $C(M)$ and $\Gamma(\mathcal{C}l M)$ (this approach is more of the spirit of noncommutative geometry). Then the notion of spin structure and finally, with the notion of spin and Clifford connection, we reach the definition of a (generalized) Dirac operator.

Here we try to bypass this approach to save time.

References: a classical book is [71], but I recommend [47]. Here, we follow [88], but see also [50].

3.1 Definition and main properties

Let (M, g) be a compact Riemannian manifold with metric g , of dimension d and E be a vector bundle over M . An example is the (Clifford) bundle $E = \mathcal{C}l T^*M$ where the fiber $\mathcal{C}l T_x^*M$ is the Clifford algebra of the real vector space T_x^*M for $x \in M$ endowed with the nondegenerate quadratic form g .

Given a connection ∇ on E , recall that a differential operator P of order m on E is an element of

$$\text{Diff}^m(M, E) := \Gamma(M, \text{End}(E)) \cdot \text{Vect}\{ \nabla_{X_1} \cdots \nabla_{X_j} \mid X_j \in \Gamma(M, TM), j \leq m \}.$$

In particular, $\text{Diff}^m(M, E)$ is a subalgebra of $\text{End}(\Gamma(M, E))$ and the operator P has a principal symbol σ_m^P in $\Gamma(T^*M, \pi^* \text{End}(E))$ where $\pi : T^*M \rightarrow M$ is the canonical submersion and $\sigma_m^P(x, \xi)$ is given by (4).

An example: Let $E = \wedge T^*M$. The exterior product and the contraction given on $\omega, \omega_j \in E$ by

$$\begin{aligned} \epsilon(\omega_1) \omega_2 &:= \omega_1 \wedge \omega_2, \\ \iota(\omega) (\omega_1 \wedge \cdots \wedge \omega_m) &:= \sum_{j=1}^m (-1)^{j-1} g(\omega, \omega_j) \omega_1 \wedge \cdots \wedge \widehat{\omega_j} \wedge \cdots \wedge \omega_m \end{aligned}$$

suggest the following definition $c(\omega) := \epsilon(\omega) + \iota(\omega)$ and one checks that

$$c(\omega_1) c(\omega_2) + c(\omega_2) c(\omega_1) = 2g(\omega_1, \omega_2) \text{id}_E. \quad (17)$$

E has a natural scalar product: if e_1, \dots, e_d is an orthonormal basis of T_x^*M , then the scalar product is chosen such that $e_{i_1} \wedge \cdots \wedge e_{i_p}$ for $i_1 < \cdots < i_p$ is an orthonormal basis.

If $d \in \text{Diff}^1$ is the exterior derivative and d^* is its adjoint for the deduced scalar product on $\Gamma(M, E)$, then their principal symbols are

$$\sigma_1^d(\omega) = i\epsilon(\omega), \quad (18)$$

$$\sigma_1^{d^*}(\omega) = -i\iota(\omega). \quad (19)$$

This follows from $\sigma_1^d(x, \xi) = \lim_{t \rightarrow \infty} \frac{1}{t} (e^{-ith(x)} d e^{ith(x)})(x) = \lim_{t \rightarrow \infty} \frac{1}{t} it d_x h = i d_x h = i \xi$ where h is such that $d_x h = \xi$, so $\sigma_1^d(x, \xi) = i \xi$ and similarly for $\sigma_1^{d^*}$.

More generally, if $P \in \text{Diff}^m(M)$, $\sigma_m^P(dh) = \frac{1}{i^m m!} (adh)^m(P)$ with $adh = [\cdot, h]$ and $\sigma_m^{P^*}(\omega) = \sigma_m^P(\omega)^*$ where the adjoint P^* is for the scalar product on $\Gamma(M, E)$ associated to an hermitean metric on E : $\langle \psi, \psi' \rangle := \int_M \langle \psi(x), \psi'(x) \rangle_x |dx|$ is a scalar product on the space $\Gamma(M, E)$.

Definition 3.1. *The operator $P \in \text{Diff}^2(M, E)$ is called a generalized Laplacian when its symbol satisfies $\sigma_2^P(x, \xi) = |\xi|_x^2 \text{id}_{E_x}$ for $x \in M$, $\xi \in T_x^*M$ (note that $|\xi|_x$ depends on the metric g).*

This is equivalent to say that, in local coordinates, $P = -\sum_{i,j} g^{ij}(x) \partial_{x^i} \partial_{x^j} + b^j(x) \partial_{x^j} + c(x)$ where the b^j are smooth and c is in $\Gamma(M, \text{End}(E))$.

Definition 3.2. *Assume that $E = E^+ \oplus E^-$ is a \mathbb{Z}_2 -graded vector bundle.*

When $D \in \text{Diff}^1(M, E)$ and $D = \begin{pmatrix} 0 & D^+ \\ D^- & 0 \end{pmatrix}$ (D is odd) where $D^\pm : \Gamma(M, E^\mp) \rightarrow \Gamma(M, E^\pm)$, D is called a Dirac operator if $D^2 = \begin{pmatrix} D^- D^+ & 0 \\ 0 & D^+ D^- \end{pmatrix}$ is a generalized Laplacian.

A good example is given by $E = \wedge T^*M = \wedge^{\text{even}} T^*M \oplus \wedge^{\text{odd}} T^*M$ and the de Rham operator $D := d + d^*$. It is a Dirac operator since $D^2 = dd^* + d^*d$ is a generalized Laplacian according to (18) (19). D^2 is also called the Laplace–Beltrami operator.

Definition 3.3. *Define $\mathcal{C}l M$ as the vector bundle over M whose fiber in $x \in M$ is the Clifford algebra $\mathcal{C}l T_x^*M$ (or $\mathcal{C}l T_x M$ using the musical isomorphism $X \in TM \leftrightarrow X^\flat \in T^*M$).*

A bundle E is called a Clifford bundle over M when there exists a \mathbb{Z}_2 -graduate action $c : \Gamma(M, \mathcal{C}l M) \rightarrow \text{End}(\Gamma(M, E))$.

The main idea which drives this definition is that Clifford actions correspond to principal symbols of Dirac operators:

Proposition 3.4. *If E is a Clifford module, every odd $D \in \text{Diff}^1$ such that $[D, f] = i c(df)$ for $f \in C^\infty(M)$ is a Dirac operator.*

Conversely, if D is a Dirac operator, there exists a Clifford action c with $c(df) = -i [D, f]$.

Proof. Let $x \in M$, $\xi \in T_x^*M$ and $f \in C^\infty(M)$ such that $d_x f = \xi$. Then

$$\sigma_1^D(df)(x) = \left(\frac{1}{i} ad f\right)D = -i[D, f] = c(df),$$

so, thanks to Theorem 1.5, $\sigma_2^{D^2}(x, \xi) = \left(\sigma_1^D(x, \xi)\right)^2 = |\xi|_x^2 \text{id}_{E_x}$ and D^2 is a generalized Laplacian.

Conversely, if D is a Dirac operator, then we can define $c(df) := i[D, f]$. This makes sense since $D \in \text{Diff}^1$ and for $f \in C^\infty(M)$, $x \in M$, $[D, f](x) = i\sigma_1^D(df)(x) = i\sigma_1^D(x, d_x f)$ is an endomorphism of E_x depending only on $d_x f$. So c can be extended to the whole T^*M with $c(x, \xi) := c(dh)(x) = i\sigma_1^D(x, \xi)$ where $h \in C^\infty(M)$ is chosen such that $\xi = d_x h$. The map $\xi \rightarrow c(x, \xi)$ is linear from T_x^*M to $\text{End}(E_x)$ and $c(x, \xi)^2 = \sigma_1^D(x, \xi)^2 = \sigma_2^{D^2}(x, \xi) = |\xi|_x^2$ for each $\xi \in T_x^*M$. Thus c can be extended to an morphism of algebras from $\mathcal{C}l(T_x^*M)$ in $\text{End}(E_x)$. This gives a Clifford action on the bundle E . \square

Consider previous example: $E = \wedge T^*M = \wedge^{\text{even}} T^*M \oplus \wedge^{\text{odd}} T^*M$ is a Clifford module for $c := i(\epsilon + \iota)$ coming from the Dirac operator $D = d + d^*$: by (18) and (19)

$$i[D, f] = i[d + d^*, f] = i\left(i\sigma_1^d(df) - i\sigma_1^{d^*}(df)\right) = -i(\epsilon + \iota)(df).$$

Definition 3.5. Let E be a Clifford module over M . A connection ∇ on E is a Clifford connection if for $a \in \Gamma(M, \mathcal{C}l M)$ and $X \in \Gamma(M, TM)$, $[\nabla_X, c(a)] = c(\nabla_X^{LC} a)$ where ∇_X^{LC} is the Levi-Civita connection after its extension to the bundle $\mathcal{C}l M$ (here $\mathcal{C}l M$ is the bundle with fiber $\mathcal{C}l T_x M$).

A Dirac operator D_∇ is associated to a Clifford connection ∇ in the following way:

$$D_\nabla := -i c \circ \nabla, \quad \Gamma(M, E) \xrightarrow{\nabla} \Gamma(M, T^*M \otimes E) \xrightarrow{c \otimes 1} \Gamma(M, E).$$

(we will use c for $c \otimes 1$).

Thus if in local coordinates, $\nabla = \sum_{j=1}^d dx^j \otimes \nabla_{\partial_j}$, the associated Dirac operator is given by $D_\nabla = -i \sum_j c(dx^j) \nabla_{\partial_j}$. In particular, for $f \in C^\infty(M)$,

$$[D_\nabla, f id_E] = -i \sum_{i=1}^d c(dx^i) [\nabla_{\partial_i}, f] = \sum_{j=1}^d -i c(dx^j) \partial_j f = -ic(df).$$

By Proposition 3.4, D_∇ deserves the name of Dirac operator!

Examples:

1) For the previous example $E = \wedge T^*M$, the Levi-Civita connection is indeed a Clifford connection whose associated Dirac operator coincides with the de Rham operator $D = d + d^*$.

2) *The spinor bundle:* Recall that the spin group Spin_d is the non-trivial two-fold covering of SO_d , so we have

$$0 \longrightarrow \mathbb{Z}_2 \longrightarrow \text{Spin}_d \xrightarrow{\xi} \text{SO}_d \longrightarrow 1.$$

Let $\text{SO}(TM) \rightarrow M$ be the SO_d -principal bundle of positively oriented orthonormal frames on TM of an oriented Riemannian manifold M of dimension d .

A *spin structure* on an oriented d -dimensional Riemannian manifold (M, g) is a Spin_d -principal bundle $\text{Spin}(TM) \xrightarrow{\pi} M$ with a two-fold covering map $\text{Spin}(TM) \xrightarrow{\eta} \text{SO}(TM)$ such that the following diagram commutes:

$$\begin{array}{ccc} \text{Spin}(TM) \times \text{Spin}_d & \longrightarrow & \text{Spin}(TM) \\ \eta \times \xi \downarrow & & \eta \downarrow \\ \text{SO}(TM) \times \text{SO}_d & \longrightarrow & \text{SO}(TM) \end{array} \begin{array}{c} \xrightarrow{\pi} \\ \xrightarrow{\pi} \end{array} M$$

where the horizontal maps are the actions of Spin_d and SO_d on the principal fiber bundles $\text{Spin}(TM)$ and $\text{SO}(TM)$.

A *spin manifold* is an oriented Riemannian manifold admitting a spin structure.

In above definition, one can replace Spin_d by the group Spin_d^c which is a central extension of SO_d by \mathbb{T} :

$$0 \longrightarrow \mathbb{T} \longrightarrow \text{Spin}_d^c \xrightarrow{\xi} \text{SO}_d \longrightarrow 1.$$

An oriented Riemannian manifold (M, g) is spin if and only if the second Stiefel–Whitney class of its tangent bundle vanishes. Thus a manifold is spin if and only both its first and second Stiefel–Whitney classes vanish (the vanishing of the first one being equivalent to the orientability of the manifold). In this case, the set of spin structures on (M, g) stands in one-to-one correspondence with $H^1(M, \mathbb{Z}_2)$. In particular the existence of a spin structure does not depend on the metric or the orientation of a given manifold. Note that all manifolds

of dimension $d \leq 4$ have spin^c structures but $\mathbb{C}P^2$ is a 4-dimensional (complex) manifold without spin structures.

Let ρ be an irreducible representation of $\mathcal{C}\ell\mathbb{C}^d \rightarrow \text{End}_{\mathbb{C}}(\Sigma_d)$ with $\Sigma_d \simeq \mathbb{C}^{2^{\lfloor d/2 \rfloor}}$ as set of complex spinors. Of course, $\mathcal{C}\ell\mathbb{C}^d$ is endowed with its canonical complex bilinear form.

The *spinor bundle* S of M is the complex vector bundle associated to the principal bundle $\text{Spin}(TM)$ with the spinor representation, namely $S := \text{Spin}(TM) \times_{\rho_d} \Sigma_d$. Here ρ_d is a representation of Spin_d on $\text{Aut}(\Sigma_d)$ which is the restriction of ρ .

More precisely, if $d = 2m$ is even, $\rho_d = \rho^+ + \rho^-$ where ρ^{\pm} are two nonequivalent irreducible complex representations of Spin_{2m} and $\Sigma_{2m} = \Sigma_{2m}^+ \oplus \Sigma_{2m}^-$, while for $d = 2m + 1$ odd, the spinor representation ρ_d is irreducible.

In practice, M is a spin manifold means that there exists a Clifford bundle $S = S^+ \oplus S^-$ such that $S \simeq \wedge T^*M$. Due to the dimension of M , the Clifford bundle has fiber

$$\mathcal{C}\ell_x M = \begin{cases} M_{2m}(\mathbb{C}) & \text{when } d = 2m \text{ is even,} \\ M_{2m}(\mathbb{C}) \oplus M_{2m}(\mathbb{C}) & \text{when } d = 2m + 1. \end{cases}$$

Locally, the spinor bundle satisfies $S \simeq M \times \mathbb{C}^{d/2}$.

A *spin connection* $\nabla^S : \Gamma^\infty(M, S) \rightarrow \Gamma^\infty(M, S) \otimes \Gamma^\infty(M, T^*M)$ is any connection which is compatible with Clifford action:

$$[\nabla^S, c(\cdot)] = c(\nabla^{LC} \cdot).$$

It is uniquely determined by the choice of a spin structure on M (once an orientation of M is chosen).

Definition 3.6. *The Dirac (also called Atiyah–Singer) operator given by the spin structure is*

$$\mathcal{D} := -i c \circ \nabla^S. \quad (20)$$

In coordinates,

$$\mathcal{D} = -i c(dx^j) (\partial_j - \omega_j(x)) \quad (21)$$

where ω_j is the spin connection part which can be computed in the coordinate basis

$$\omega_j = \frac{1}{4} \left(\Gamma_{ji}^k g_{kl} - \partial_i (h_j^\alpha) \delta_{\alpha\beta} h_i^\beta \right) c(dx^i) c(dx^l)$$

where the matrix $H := [h_j^\alpha]$ is such that $H^t H = [g_{ij}]$ (we use Latin letters for coordinate basis indices and Greek letters for orthonormal basis indices).

This gives $\sigma_1^D(x, \xi) = c(\xi) + i c(dx^j) \omega_j(x)$. Thus in *normal coordinates around x_0* ,

$$\begin{aligned} c(dx^j)(x_0) &= \gamma^j, \\ \sigma_1^D(x_0, \xi) &= c(\xi) = \gamma^j \xi_j \end{aligned}$$

where the γ 's are constant hermitean matrices.

A fundamental result concerning a Dirac operator (definition 3.2) is its unique continuation property: if ψ satisfies $D\psi = 0$ and ψ vanishes on an open subset of the smooth manifold

M (with or without boundary), then ψ also vanishes on the whole connected component of M .

The Hilbert space of spinors is

$$\mathcal{H} = L^2((M, g), S) := \{ \psi \in \Gamma^\infty(M, S) \mid \int_M \langle \psi, \psi \rangle_x d\text{vol}_g(x) < \infty \} \quad (22)$$

where we have a scalar product which is $C^\infty(M)$ -valued. On its domain $\Gamma^\infty(M, S)$, the Dirac operator is symmetric: $\langle \psi, \mathcal{D}\phi \rangle = \langle \mathcal{D}\psi, \phi \rangle$. Moreover, it has a selfadjoint closure (which is \mathcal{D}^{**}):

Theorem 3.7. *Let (M, g) be an oriented compact Riemannian spin manifold without boundary. By extension to \mathcal{H} , \mathcal{D} is essentially selfadjoint on its original domain $\Gamma^\infty(M, S)$. It is a differential (unbounded) operator of order one which is elliptic.*

See [50, 71, 106, 107] for a proof.

There is a nice formula which relates the Dirac operator \mathcal{D} to the spinor Laplacian

$$\Delta^S := -\text{Tr}_g(\nabla^S \circ \nabla^S) : \Gamma^\infty(M, S) \rightarrow \Gamma^\infty(M, S).$$

Before to give it, we need to fix few notations: let $R \in \Gamma^\infty(M, \wedge^2 T^*M \otimes \text{End}(TM))$ be the *Riemann curvature tensor* with components $R_{ijkl} := g(\partial_i, R(\partial_k, \partial_l)\partial_j)$, the *Ricci tensor* components are $R_{jl} := g^{ik}T_{ijkl}$ and the *scalar curvature* is $s := g^{jl}R_{jl}$.

Proposition 3.8. *Schrödinger–Lichnerowicz formula: with same hypothesis,*

$$\mathcal{D}^2 = \Delta^S + \frac{1}{4}s \quad (23)$$

where s is the scalar curvature of M .

The proof is just a lengthy computation (see for instance [50]).

We already know via Theorems 1.10 and 3.7 that \mathcal{D}^{-1} is compact so has a discrete spectrum. For $T \in \mathcal{K}_+(\mathcal{H})$, we denote by $\{\lambda_n(T)\}_{n \in \mathbb{N}}$ its spectrum sorted in decreasing order including multiplicity (and in increasing order for an unbounded positive operator T such that T^{-1} is compact) and by $N_T(\lambda) := \#\{\lambda_n(T) \mid \lambda_n \leq \lambda\}$ its counting function.

Theorem 3.9. *With same hypothesis, the asymptotics of the Dirac operator counting function is $N_{|\mathcal{D}|}(\lambda) \underset{\lambda \rightarrow \infty}{\sim} \frac{2^d \text{Vol}(\mathbb{S}^{d-1})}{d(2\pi)^d} \text{Vol}(M) \lambda^d$ where $\text{Vol}(M) = \int_M d\text{vol}$.*

Proof. By Weyl's theorem, we know the asymptotics of $N_\Delta(\lambda)$ for the the scalar Laplacian $\Delta := -\text{Tr}_g(\nabla^{T^*M \otimes T^*M} \circ \nabla^{T^*M})$ which in coordinates is $\Delta = -g^{ij}(\partial_i \partial_j - \Gamma_{ij}^k \partial_k)$. It is given by:

$$N_\Delta(\lambda) \underset{\lambda \rightarrow \infty}{\sim} \frac{\text{Vol}(\mathbb{S}^{d-1})}{d(2\pi)^d} \text{Vol}(M) \lambda^{d/2}.$$

For the spinor Laplacian, we get the same formula with an extra factor of $\text{Tr}(1_S) = 2^d$ and Proposition 3.8 shows that $N_{\mathcal{D}^2}(\lambda)$ has the same asymptotics than $N_\Delta(\lambda)$ since s gives rise to a bounded operator. \square

We already encounter such computation in Example 2.13.

3.2 Dirac operators and change of metrics

Recall that the spinor bundle S_g and square integrable spinors \mathcal{H}_g defined in (22) depends on the chosen metric g , so we note M_g instead of M and $\mathcal{H}_g := L^2(M_g, S_g)$ and a natural question is: what happens to a Dirac operator when the metric changes?

Let g' be another Riemannian metric on M . Since the space of d -forms is one-dimensional, there exists a positive function $f_{g,g'} : M \rightarrow \mathbb{R}^+$ such that $dvol_{g'} = f_{g',g} dvol_g$.

Let $I_{g,g'}(x) : S_g \rightarrow S_{g'}$ the natural injection on the spinors spaces above point $x \in M$ which is a pointwise linear isometry: $|I_{g,g'}(x) \psi(x)|_{g'} = |\psi(x)|_g$. Let us first see its construction: there always exists a g -symmetric automorphism $H_{g,g'}$ of the $2^{\lfloor d/2 \rfloor}$ -dimensional vector space TM such that $g'(X, Y) = g(H_{g,g'} X, Y)$ for $X, Y \in TM$ so define $\iota_{g,g'} X := H_{g,g'}^{-1/2} X$. Note that $\iota_{g,g'}$ commutes with right action of the orthogonal group O_d and can be lifted up to a diffeomorphism Pin_d -equivariant on the spin structures associated to g and g' and this lift is denoted by $I_{g,g'}$ (see [6]). This isometry is extended as operator on the Hilbert spaces $I_{g,g'} : \mathcal{H}_g \rightarrow \mathcal{H}_{g'}$ with $(I_{g,g'} \psi)(x) := I_{g,g'}(x) \psi(x)$.

Now define

$$U_{g',g} := \sqrt{f_{g,g'}} I_{g',g} : \mathcal{H}_{g'} \rightarrow \mathcal{H}_g. \quad (24)$$

Then by construction, $U_{g',g}$ is a unitary operator from $\mathcal{H}_{g'}$ onto \mathcal{H}_g : for $\psi' \in \mathcal{H}_{g'}$,

$$\begin{aligned} \langle U_{g',g} \psi', U_{g',g} \psi' \rangle_{\mathcal{H}_g} &= \int_M |U_{g',g} \psi'|_g^2 dvol_g = \int_M |I_{g',g} \psi'|_g^2 f_{g',g} dvol_g = \int_M |\psi'|_{g'}^2 dvol_{g'} \\ &= \langle \psi', \psi' \rangle_{\mathcal{H}_{g'}}. \end{aligned}$$

So we can realize $\mathcal{D}_{g'}$ as an operator $D_{g'}$ acting on \mathcal{H}_g with

$$D_{g'} : \mathcal{H}_g \rightarrow \mathcal{H}_g, \quad D_{g'} := U_{g',g}^{-1} \mathcal{D}_{g'} U_{g',g}. \quad (25)$$

This is an unbounded operator on \mathcal{H}_g which has the same eigenvalues as $\mathcal{D}_{g'}$.

In the same vein, the k -th Sobolev space $H^k(M_g, S_g)$ (which is the completion of the space $\Gamma^\infty(M_g, S_g)$ under the norm $\|\psi\|_k^2 = \sum_{j=0}^k \int_M |\nabla^j \psi(x)|^2 dx$; be careful, ∇ applied to $\nabla \psi$ is the tensor product connection on $T^*M_g \otimes S_g$ etc, see Theorem 1.10) can be transported: the map $U_{g,g'} : H^k(M_g, S_g) \rightarrow H^k(M_{g'}, S_{g'})$ is an isomorphism, see [98]. In particular, (after the transport map U), the domain of $D_{g'}$ and $\mathcal{D}_{g'}$ are the same.

A nice example of this situation is when g' is in the conformal class of g where we can compute explicitly $\mathcal{D}_{g'}$ and $D_{g'}$ [2, 6, 47, 57].

Theorem 3.10. *Let $g' = e^{2h}g$ be a conformal transformation of g with $h \in C^\infty(M, \mathbb{R})$. Then there exists an isometry $I_{g,g'}$ between the spinor bundle S_g and $S_{g'}$ such that*

$$\begin{aligned} \mathcal{D}_{g'} I_{g,g'} \psi &= e^{-h} I_{g,g'} \left(\mathcal{D}_g \psi - i \frac{d-1}{2} c_g(\text{grad } h) \psi \right), \\ \mathcal{D}_{g'} &= e^{-\frac{d+1}{2}h} I_{g,g'} \mathcal{D}_g I_{g,g'}^{-1} e^{\frac{d-1}{2}h}, \\ D_{g'} &= e^{-h/2} \mathcal{D}_g e^{-h/2}. \end{aligned}$$

for $\psi \in \Gamma^\infty(M, S_g)$.

Proof. The isometry $X \rightarrow X' := e^{-h}X$ from (TM, g) onto (TM, g') defines a principal bundle isomorphism $SO_g(TM) \rightarrow SO_{g'}(TM)$ lifting to the spin level. More precisely, it induces a vector-bundle isomorphism $I_{g,g'} : S_g \rightarrow S_{g'}$, preserving the pointwise hermitean inner product (i.e. $H_{g,g'} = e^{2h}$), such that $e^{-h}c_{g'}(X)I_{g,g'}\psi = c_{g'}(X')I_{g,g'}\psi = I_{g,g'}c_g(X)\psi$.

For a connection $\widetilde{\nabla}$ compatible with a metric k and without torsion, we have for X, Y, Z in $\Gamma^\infty(TM)$

$$\begin{aligned} 2k(\widetilde{\nabla}_X Y, Z) &= k([X, Y], Z) + k([Z, Y], X) + k([Z, X], Y) \\ &\quad + X \cdot k(Y, Z) + Y \cdot k(X, Z) - Z \cdot k(X, Y) \end{aligned} \quad (26)$$

which is obtained via $k(\widetilde{\nabla}_X Y, Z) + k(Y, \widetilde{\nabla}_X Z) = X \cdot k(Y, Z)$ minus two cyclic permutations. The set $\{e'_j := e^{-h}e_j \mid 1 \leq j \leq d\}$ is a local g' -orthonormal basis of TU for g' if and only if $\{e_j\}$ is a local g -orthonormal basis of TU where U is a trivializing open subset of M . Applying (26) to ∇'^{LC} , we get

$$\begin{aligned} 2g'(\nabla'^{LC} e'_i, e'_j) &= e^{2h}g([X, e^{-h}e_i], e^{-h}e_j) + e^{2h}g([e^{-h}e_j, e^{-h}e_i], X) + e^{2h}g([e^{-h}e_j, X], e^{-h}e_i) \\ &\quad + X \cdot e^{2h}g(e^{-h}e_i, e^{-h}e_j) + e^{-h}e_i \cdot e^{2h}g(X, e^{-h}e_j) - e^{-h}e_j \cdot e^{2h}g(X, e^{-h}e_i) \\ &= 2g(\nabla_X^{LC} e_i, e_j) + 2(e_i \cdot h)g(e_j, X) - 2(e_j \cdot h)g(e_i, X). \end{aligned}$$

Since $\nabla_X^{S_g} \psi = -\frac{1}{4}g(\nabla_X^{LC} e_i, e_j) c_g(e_i) c_g(e_j) \psi$, for $\psi \in \Gamma^\infty(U, S_g)$,

$$\nabla_{e'_k}^{S_{g'}} I_{g,g'} \psi = I_{g,g'} \left[\nabla_{e_k}^{S_g} + \frac{1}{2}c_g(e_k) c_g(\text{grad } h) - \frac{1}{2}e_k(h) \right] \psi. \quad (27)$$

Hence

$$\begin{aligned} \mathbb{D}_{g'} I_{g,g'} \psi &= -ic_{g'}(e'_k) \nabla_{e'_k}^{S_{g'}} I_{g,g'} \psi = -ie^{-h} c_{g'}(e'_k) \nabla_{e_k}^{S_{g'}} I_{g,g'} \psi \\ &= -ie^{-h} c_{g'}(e'_k) I_{g,g'} \left[\nabla_{e_k}^{S_g} + \frac{1}{2}c_g(e_k) c_g(\text{grad } h) - \frac{1}{2}e_k(h) \right] \psi \\ &= -ie^{-h} I_{g,g'} c_g(e_k) \left[\nabla_{e_k}^{S_g} + \frac{1}{2}c_g(e_k) c_g(\text{grad } h) - \frac{1}{2}e_k(h) \right] \psi \\ &= e^{-h} I_{g,g'} \left[\mathbb{D}_g \psi - i\frac{d-1}{2} c_g(\text{grad } h) \right] \psi. \end{aligned}$$

So, using $[\mathbb{D}_g, f] = -ic_g(\text{grad } f)$ for $f = e^{-\frac{d-1}{2}h}$,

$$\begin{aligned} \mathbb{D}_{g'} e^{-\frac{d-1}{2}h} I_{g,g'} \psi &= e^{-h} I_{g,g'} \left[\mathbb{D}_g e^{-\frac{d-1}{2}h} \psi - i\frac{d-1}{2} e^{-\frac{d-1}{2}h} c_g(\text{grad } h) \psi \right] \\ &= e^{-h} I_{g,g'} \left[e^{-\frac{d-1}{2}h} \mathbb{D}_g \psi + [\mathbb{D}_g, e^{-\frac{d-1}{2}h}] \psi - i\frac{d-1}{2} e^{-\frac{d-1}{2}h} c_g(\text{grad } h) \psi \right] \\ &= e^{-h} I_{g,g'} \left[e^{-\frac{d-1}{2}h} \mathbb{D}_g \psi - \frac{d-1}{2} e^{-\frac{d-1}{2}h} (-i)c_g(\text{grad } h) \psi \right. \\ &\quad \left. - i\frac{d-1}{2} e^{-\frac{d-1}{2}h} c_g(\text{grad } h) \psi \right] \\ &= e^{-\frac{d+1}{2}h} I_{g,g'} \mathbb{D}_g \psi. \end{aligned}$$

Thus $\mathbb{D}_{g'} = e^{-\frac{d+1}{2}h} I_{g,g'} \mathbb{D}_g I_{g,g'}^{-1} e^{\frac{d-1}{2}h}$ and since $dvol_{g'} = e^{dh} dvol_g$, using (24)

$$\mathbb{D}_{g'} = e^{-h/2} U_{g',g}^{-1} \mathbb{D}_g U_{g',g} e^{-h/2}.$$

Finally, (25) yields $D_{g'} = e^{-h/2} \mathbb{D}_g e^{-h/2}$. □

Note that $D_{g'}$ is not a Dirac operator as defined in (20) since its principal symbol has an x -dependence: $\sigma^{D_{g'}}(x, \xi) = e^{-h(x)} c_g(\xi)$.

The principal symbols of $\mathcal{D}_{g'}$ and \mathcal{D}_g are related by

$$\sigma_d^{\mathcal{D}_{g'}}(x, \xi) = e^{-h(x)/2} U_{g',g}^{-1}(x) \sigma_d^{\mathcal{D}_g}(x, \xi) U_{g',g}(x) e^{-h(x)/2}, \quad \xi \in T_x^* M.$$

Thus

$$c_{g'}(\xi) = e^{-h(x)} U_{g',g}^{-1}(x) c_g(\xi) U_{g',g}(x), \quad \xi \in T_x^* M. \quad (28)$$

Using $c_g(\xi)c_g(\eta) + c_g(\eta)c_g(\xi) = 2g(\xi, \eta) \text{id}_{S_g}$, formula (28) gives a verification of the formula $g'(\xi, \eta) = e^{-2h} g(\xi, \eta)$.

Note that two volume forms μ, μ' on a compact connected manifold M are related by an orientation preserving diffeomorphism α of M in the following sense [81]: there exists a constant $c = (\int_M \mu')^{-1} \int_M \mu$ such that $\mu = c \alpha^* \mu'$ where $\alpha^* \mu'$ is the pull-back of μ' (i.e. $\int_{\alpha(S)} \alpha^* \mu = \int_S \mu$ for any set $S \subset M$). The proof is based on the construction of an orientation preserving automorphism homotopic to the identity.

It is also natural to look at the changes on a Dirac operator when the metric g is modified by a diffeomorphism α which preserves the spin structure. The diffeomorphism α can be lifted to a diffeomorphism O_d -equivariant on the O_d -principal bundle of g -orthonormal frames with $\tilde{\alpha} := H_{\alpha^*g}^{-1/2} T\alpha$, and this lift also exists on S_g when α preserves both the orientation and the spin structure. However, the last lift is defined up to a \mathbb{Z}_2 -action which disappears if α is connected to the identity.

The pull-back $g' := \alpha^* g$ of the metric g is defined by $(\alpha^* g)_x(\xi, \eta) = g_{\alpha(x)}(\alpha_*(\xi), \alpha_*(\eta))$, $x \in M$, where α_* is the push-forward map: $T_x M \rightarrow T_{\alpha(x)} M$. Of course, the metric g' and g are different but the geodesic distances are the same. Let us check that $d_{g'} = \alpha^* d_g$:

In local coordinates, we note $\partial_\mu := \partial/\partial x^\mu$ and $\partial'_{\mu'} := \partial/\partial(\alpha(x))^{\mu'}$. Thus $\partial' = (\Lambda^{-1T}) \partial$ where $\Lambda^{\mu'}{}_\mu := \partial(\alpha(x))^{\mu'}/\partial x^\mu$. The dependence in the metric g of Cristoffel symbols is $\Gamma_{\mu\nu}^\rho = \frac{1}{2} g^{\rho\beta} (\partial_\mu g_{\beta\nu} + \partial_\nu g_{\mu\beta} - \partial_\beta g_{\mu\nu})$. Thus the same symbols Γ' associated to g' are

$$\Gamma'^{\rho'}{}_{\mu'\nu'} = \Lambda^{\rho'}{}_\rho (\Lambda^{-1T})_{\mu'}{}^\mu (\Lambda^{-1T})_{\nu'}{}^\nu \Gamma_{\mu\nu}^\rho + \Lambda^{\rho'}{}_\rho (\Lambda^{-1T})_{\mu'}{}^\mu \partial_\mu (\Lambda^{-1T})_{\nu'}{}^\nu. \quad (29)$$

The geodesic equation is $\ddot{x}^\rho + \Gamma_{\mu\nu}^\rho \dot{x}^\mu \dot{x}^\nu = 0$ for all ρ (note that neither x^μ nor \dot{x}^μ are 4-vectors in the sense that they are not transformed like $v^{\mu'} = \Lambda^{\mu'}{}_\mu v^\mu$, while \dot{x}^μ is a 4-vector; in fact $\alpha(\ddot{x})^{\rho'} = \Lambda^{\rho'}{}_\rho \ddot{x}^\rho + \partial_\mu \Lambda^{\rho'}{}_\rho \dot{x}^\mu \dot{x}^\rho$). This relation and (29) give the invariance of the geodesic equation and the same for the distance since for any path γ joining points $x = \gamma(0)$, $y = \gamma(1)$

$$\int_0^1 \sqrt{(\alpha^* g)_{\gamma(t)}(\gamma'(t), \gamma'(t))} dt = \int_0^1 \sqrt{g_{\alpha \circ \gamma(t)}((\alpha \circ \gamma)'(t), (\alpha \circ \gamma)'(t))} dt$$

and $(\alpha \circ \gamma)(0) = \alpha(x)$, $(\alpha \circ \gamma)(1) = \alpha(y)$. Note that α is an isometry only if $\alpha^* d_g = d_g$.

Recall that the principal symbol of a Dirac operator D is $\sigma_d^D(x, \xi) = c_g(\xi)$ so gives the metric g by (17) as we checked above. This information will be used later in the definition of a spectral triple. A commutative spectral triple associated to a manifold generates the so-called Connes' distance which is nothing else but the metric distance; see the remark after (43). Again, the link between d_{α^*g} and d_g is explained by (25), since the unitary induces an automorphism of the C^* -algebra $C^\infty(M)$.

4 Heat kernel expansion

References for this section: [4, 44, 45] and especially [109].

Recall that the heat kernel is a Green function of the heat operator $e^{t\Delta}$ (recall that $-\Delta$ is a positive operator) which measures the temperature evolution in a domain whose boundary has a given temperature. For instance, the heat kernel of the Euclidean space \mathbb{R}^d is

$$k_t(x, y) = \frac{1}{(4\pi t)^{d/2}} e^{-|x-y|^2/4t} \text{ for } x \neq y \quad (30)$$

and it solves the heat equation

$$\begin{cases} \partial_t k_t(x, y) = \Delta_x k_t(x, y), & \forall t > 0, x, y \in \mathbb{R}^d \\ \text{initial condition: } \lim_{t \downarrow 0} k_t(x, y) = \delta(x - y). \end{cases}$$

Actually, $k_t(x, y) = \frac{1}{2\pi} \int_{-\infty}^{\infty} e^{-ts^2} e^{is(x-y)} ds$ when $d = 1$.

Note that for $f \in \mathcal{D}(\mathbb{R}^d)$, we have $\lim_{t \downarrow 0} \int_{\mathbb{R}^d} k_t(x, y) f(y) dy = f(x)$.

For a connected domain (or manifold with boundary with vector bundle V) U , let λ_n be the eigenvalues for the Dirichlet problem of minus the Laplacian

$$\begin{cases} -\Delta\phi = \lambda\psi & \text{in } U \\ \psi = 0 & \text{on } \partial U. \end{cases}$$

If $\psi_n \in L^2(U)$ are the normalized eigenfunctions, the inverse Dirichlet Laplacian Δ^{-1} is a selfadjoint compact operator, $0 \leq \lambda_1 \leq \lambda_2 \leq \lambda_3 \leq \dots, \lambda_n \rightarrow \infty$.

The interest for the heat kernel is that, if $f(x) = \int_0^\infty dt e^{-tx} \phi(x)$ is the Laplace transform of ϕ , then $\text{Tr}(f(-\Delta)) = \int_0^\infty dt \phi(t) \text{Tr}(e^{t\Delta})$ (if everything makes sense) is controlled by $\text{Tr}(e^{t\Delta}) = \int_M d\text{vol}(x) \text{tr}_{V_x} k_t(x, x)$ since $\text{Tr}(e^{t\Delta}) = \sum_{n=1}^\infty e^{t\lambda_n}$ and

$$k_t(x, y) = \langle x, e^{t\Delta} y \rangle = \sum_{n,m=1}^\infty \langle x, \psi_m \rangle \langle \psi_m, e^{t\Delta} \psi_n \rangle \langle \psi_n, y \rangle = \sum_{n,m=1}^\infty \overline{\psi_n(x)} \psi_n(y) e^{t\lambda_n}.$$

So it is useful to know the asymptotics of the heat kernel k_t on the diagonal of $M \times M$ especially near $t = 0$.

4.1 The asymptotics of heat kernel

Let now M be a smooth compact Riemannian manifold without boundary, V be a vector bundle over M and $P \in \Psi DO^m(M, V)$ be a positive elliptic operator of order $m > 0$. If $k_t(x, y)$ is the kernel of the heat operator e^{-tP} , then the following asymptotics exists on the diagonal:

$$k_t(x, x) \underset{t \downarrow 0^+}{\sim} \sum_{k=0}^\infty a_k(x) t^{(-d+k)/m}$$

which means that

$$\left| k_t(x, x) - \sum_{k \leq k(n)} a_k(x) t^{(-d+k)/m} \right|_{\infty, n} < c_n t^n \text{ for } 0 < t < 1$$

where $|f|_{\infty,n} := \sup_{x \in M} \sum_{|\alpha| \leq n} |\partial_x^\alpha f|$ (since P is elliptic, $k_t(x, y)$ is a smooth function of (t, x, y) for $t > 0$, see [44, section 1.6, 1.7]).

More generally, we will use

$$k(t, f, P) := \text{Tr} \left(f e^{-tP} \right)$$

where f is a smooth function. We have similarly

$$k(t, f, P) \underset{t \downarrow 0^+}{\sim} \sum_{k=0}^{\infty} a_k(f, P) t^{(-d+k)/m}. \quad (31)$$

The utility of function f will appear later for the computation of coefficients a_k . The following points are of importance:

1) The existence of this asymptotics is non-trivial [44, 45].

2) The coefficients $a_{2k}(f, P)$ can be computed locally as integral of local invariants: Recall that a locally computable quantity is the integral on the manifold of a local frame-independent smooth function of one variable, depending only on a finite number of derivatives of a finite number of terms in the asymptotic expansion of the total symbol of P . In noncommutative geometry, local generally means that it is concentrated at infinity in momentum space.

3) The odd coefficients are zero: $a_{2k+1}(f, P) = 0$.

For instance, let us assume from now on that P is a Laplace type operator of the form

$$P = -(g^{\mu\nu} \partial_\mu \partial_\nu + \mathbb{A}^\mu \partial_\mu + \mathbb{B}) \quad (32)$$

where $(g^{\mu\nu})_{1 \leq \mu, \nu \leq d}$ is the inverse matrix associated to the metric g on M , and \mathbb{A}^μ and \mathbb{B} are smooth $L(V)$ -sections on M (endomorphisms) (see also Definition 3.1). Then (see [45, Lemma 1.2.1]) there is a unique connection ∇ on V and a unique endomorphism E such that

$$P = -(\text{Tr}_g \nabla^2 + E), \quad \nabla^2(X, Y) := [\nabla_X, \nabla_Y] - \nabla_{\nabla_X^{LC} Y},$$

X, Y are vector fields on M and ∇^{LC} is the Levi-Civita connection on M . Locally

$$\text{Tr}_g \nabla^2 := g^{\mu\nu} (\nabla_\mu \nabla_\nu - \Gamma_{\mu\nu}^\rho \nabla_\rho)$$

where $\Gamma_{\mu\nu}^\rho$ are the Christoffel coefficients of ∇^{LC} . Moreover (with local frames of T^*M and V), $\nabla = dx^\mu \otimes (\partial_\mu + \omega_\mu)$ and E are related to $g^{\mu\nu}$, \mathbb{A}^μ and \mathbb{B} through

$$\omega_\nu = \frac{1}{2} g_{\nu\mu} (\mathbb{A}^\mu + g^{\sigma\varepsilon} \Gamma_{\sigma\varepsilon}^\mu \text{id}_V), \quad (33)$$

$$E = \mathbb{B} - g^{\nu\mu} (\partial_\nu \omega_\mu + \omega_\nu \omega_\mu - \omega_\sigma \Gamma_{\nu\mu}^\sigma). \quad (34)$$

In this case, the coefficients $a_k(f, P) = \int_M d\text{vol}_g \text{tr}_E (f(x) a_k(P)(x))$ and the $a_k(P) = c_i \alpha_k^i(P)$ are linear combination with constants c_i of all possible independent invariants $\alpha_k^i(P)$ of dimension k constructed from E, Ω, R and their derivatives (Ω is the curvature of the connection ω , and R is the Riemann curvature tensor). As an example, for $k = 2$, E and s are the only independent invariants.

Point 3) follow since there is no odd-dimension invariant.

4.2 Computations of heat kernel coefficients

The computation of coefficients $a_k(f, P)$ is made by induction using first a variational method: for any smooth functions f, h on has

$$\frac{d}{d\epsilon}\Big|_{\epsilon=0} a_k(1, e^{-2\epsilon f} P) = (d - k) a_k(f, P), \quad (35)$$

$$\frac{d}{d\epsilon}\Big|_{\epsilon=0} a_k(1, P - \epsilon h) = a_{k-2}(h, P), \quad (36)$$

$$\frac{d}{d\epsilon}\Big|_{\epsilon=0} a_{d-2}(e^{-2\epsilon f} h, e^{-2\epsilon f} P) = 0. \quad (37)$$

The first equation follows from

$$\frac{d}{d\epsilon}\Big|_{\epsilon=0} \text{Tr} \left(e^{-e^{-2\epsilon f} tP} \right) = -2t \frac{d}{dt} \text{Tr} \left(f e^{-tP} \right)$$

with an expansion in power series in t . Same method for (36).

For the proof of (37), we use $P(\epsilon, \delta) := e^{-2f}(P - \delta h)$; with (35) for $k = d$,

$$0 = \frac{d}{d\epsilon}\Big|_{\epsilon=0} a_d(1, P(\epsilon, \delta)),$$

thus after a variation of δ ,

$$0 = \frac{d}{d\delta}\Big|_{\delta=0} \frac{d}{d\epsilon}\Big|_{\epsilon=0} a_d(1, P(\epsilon, \delta)) = \frac{d}{d\epsilon}\Big|_{\epsilon=0} \frac{d}{d\delta}\Big|_{\delta=0} a_d(1, P(\epsilon, \delta)),$$

we derive (37) from (36).

The idea behind equations (35), (36) and (37) is that (36) shows dependence of coefficients a_k on E , while the two others describe their behaviors under local scale transformations.

Then, the $a_k(P) = c_i \alpha_k^i(P)$ are computed with arbitrary constants c_i (they are dependent only of the dimension d) and these constants are inductively calculated using (35), (36) and (37). If s is the scalar curvature and ‘;’ denote multiple covariant derivative with respect to Levi-Civita connection on M , one finds, with rescaled α 's,

$$\begin{aligned} a_0(f, P) &= (4\pi)^{-d/2} \int_M d\text{vol}_g \text{tr}_V(\alpha_0 f), \\ a_2(f, P) &= \frac{(4\pi)^{-d/2}}{6} \int_M d\text{vol}_g \text{tr}_V[(f(\alpha_1 E + \alpha_2 s))], \\ a_4(f, P) &= \frac{(4\pi)^{-d/2}}{360} \int_M d\text{vol}_g \text{tr}_V \left[f(\alpha_3 E_{;kk} + \alpha_4 E s + \alpha_5 E^2 + \alpha_6 R_{;kk} + \alpha_7 s^2 \right. \\ &\quad \left. + \alpha_8 R_{ij} R_{ij} + \alpha_9 R_{ijkl} R_{ijkl} + \alpha_{10} \Omega_{ij} \Omega_{ij}) \right]. \end{aligned} \quad (38)$$

In a_4 , they are no other invariants: for instance, $R_{ij;ij}$ is proportional to $R_{;ij}$.

Using the scalar Laplacian on the circle, one finds $\alpha_0 = 1$.

Using (36) with $k = 2$, under the change $P \rightarrow P - \epsilon h$, E becomes $E + \epsilon h$, so

$$\frac{1}{6} \int_M d\text{vol}_g \text{tr}_V(\alpha_1 h) = \int_M d\text{vol}_g \text{tr}_V(h)$$

yielding $\alpha_1 = 6$. For $k = 4$, it gives now:

$$\frac{1}{360} \int_M d\text{vol}_g \text{tr}_V(\alpha_4 h s + 2\alpha_5 h E) = \frac{1}{6} \int_M d\text{vol}_g \text{tr}_V(\alpha_1 h E + \alpha_2 h s),$$

thus $\alpha_5 = 180$ and $\alpha_4 = 60\alpha_2$.

To go further, one considers the scale transformation on P given in (35) and (37). In (35), P is transformed covariantly, the metric g is changed into $e^{-2\epsilon f}g$ implying conformal transformation of the Riemann tensor, Ricci tensor and scalar curvature giving the modifications on ω and E via (33), (34). This gives (we collect here all terms appearing in a_2 and only few terms appearing in a_4)

$$\begin{aligned}
\frac{d}{d\epsilon}|_{\epsilon=0} d\text{vol}_g &= d f d\text{vol}_g, \\
\frac{d}{d\epsilon}|_{\epsilon=0} E &= -2fE + \frac{1}{2}(d-2) f_{;ii}, \\
\frac{d}{d\epsilon}|_{\epsilon=0} s &= -2fs - 2(d-1) f_{;ii}, \\
\frac{d}{d\epsilon}|_{\epsilon=0} Es &= -4fEs + \frac{1}{2}(d-2)s f_{;ii} - 2(d-1) f_{;ii}E, \\
\frac{d}{d\epsilon}|_{\epsilon=0} E^2 &= -4fE^2 + (d-2) f_{;ii}E, \\
\frac{d}{d\epsilon}|_{\epsilon=0} s^2 &= -4fs^2 - 4(d-1) f_{;ii}s, \\
\frac{d}{d\epsilon}|_{\epsilon=0} R_{ijkl} &= -2fR_{ijkl} + \delta_{jl}f_{;ik} + \delta_{ik}f_{;jl} - \delta_{il}f_{;jk} - \delta_{jk}f_{;il}, \\
\frac{d}{d\epsilon}|_{\epsilon=0} \Omega_{ij}\Omega_{ij} &= -4f \Omega_{ij}\Omega_{ij}, \\
&\dots
\end{aligned}$$

Applying (37) with $d = 4$, we get

$$\frac{d}{d\epsilon}|_{\epsilon=0} a_2(e^{-2\epsilon f}h, e^{-2\epsilon f}P) = 0.$$

Picking terms with $\int_M d\text{vol}_g \text{tr}_V(hf_{;ii})$, we find $\alpha_1 = 6\alpha_2$, so $\alpha_2 = 1$ and $\alpha_4 = 60$. Thus $a_2(f, P)$ has been determined.

Similar method gives $a_4(f, P)$, but only after lengthy computation despite the use of Gauss–Bonnet theorem for the determination of α_{10} ! One finds:

$$\alpha_3 = 60, \alpha_5 = 180, \alpha_6 = 12, \alpha_7 = 5, \alpha_8 = -2, \alpha_9 = 2, \alpha_{10} = 30.$$

The coefficient a_6 was computed by Gilkey, a_8 by Amsterdamski, Berkin and O’Connor and a_{10} in 1998 by van de Ven [110]. Some higher coefficients are known in flat spaces.

4.3 Wodzicki residue and heat expansion

Wodzicki has proved that, in (31), $a_k(P)(x) = \frac{1}{m} c_{P^{(k-d)/m}}(x)$ is true not only for $k = 0$ as seen in Theorem 2.14 (where $P \leftrightarrow P^{-1}$), but for all $k \in \mathbb{N}$. In this section, we will prove this result when P is the inverse of a Dirac operator and this will be generalized in the next section.

Let M be a compact Riemannian manifold of dimension d even, E a Clifford module over M and D be the Dirac operator (definition 3.2) given by a Clifford connection on E . By Theorem 3.7, D is a selfadjoint (unbounded) operator on $\mathcal{H} := L^2(M, S)$.

We are going to use the heat operator e^{-tD^2} since D^2 is related to the Laplacian via the Schrödinger–Lichnerowicz formula (3.2) and since the asymptotics of the heat kernel of this Laplacian is known.

For $t > 0$, we have $e^{-tD^2} \in \mathcal{L}^1$: the result follows from the decomposition

$$e^{-tD^2} = (1 + D^2)^{(d+1)/2} e^{-tD^2} (1 + D^2)^{-(d+1)/2},$$

since $(1 + D^2)^{-(d+1)/2} \in \mathcal{L}^1$ and the function: $\lambda \rightarrow (1 + \lambda^2)^{(d+1)/2} e^{-t\lambda^2}$ is bounded.

Thus $\text{Tr}(e^{-tD^2}) = \sum_n e^{-t\lambda_n^2} < \infty$.

Another argument is the following: $(1 + D^2)^{-d/2}$ maps $L^2(M, S)$ into the Sobolev space $H^k(M, S)$ (see Theorem 1.10) and the injection $H^k(M, S) \hookrightarrow L^2(M, S)$ is Hilbert–Schmidt operator for $k > \frac{1}{2}d$. Thus $t \rightarrow e^{-tD^2}$ is a semigroup of Hilbert–Schmidt operators for $t > 0$.

Moreover, the operator e^{-tD^2} has a smooth kernel since it is regularizing, see Remark 1.3 (or [71]) and the asymptotics of its kernel is (recall 30):

$$k_t(x, y) \underset{t \downarrow 0^+}{\sim} \frac{1}{(4\pi t)^{d/2}} \sqrt{\det g_x} \sum_{j \geq 0} k_j(x, y) t^j e^{-d_g(x, y)^2/4t}$$

where k_j is a smooth section on $E^* \otimes E$. Thus

$$\text{Tr}(e^{-tD^2}) \underset{t \downarrow 0^+}{\sim} \sum_{j \geq 0} t^{(j-d)/2} a_j(D^2) \tag{39}$$

with for $j \in \mathbb{N}$,

$$\begin{cases} a_{2j}(D^2) := \frac{1}{(4\pi)^{d/2}} \int_M \text{tr}(k_j(x, x)) \sqrt{\det g_x} |dx|, \\ a_{2j+1}(D^2) = 0. \end{cases}$$

The aim now is to compute $W\text{Res}(D^{-p})$ for $0 \leq p \leq d$:

Theorem 4.1. $D^{-p} \in \Psi DO^{-p}(M, E)$ and

$$W\text{Res}(D^{-p}) = \frac{2}{\Gamma(p/2)} a_{d-p}(D^2) = \frac{2}{(4\pi)^{d/2} \Gamma(p/2)} \int_M \text{tr}(k_{(d-p)/2}(x, x)) \, d\text{vol}_g(x).$$

Proof. Assume D is invertible, otherwise swap D for the invertible operator $D + P$ where P is the projection on the kernel of D . Since the kernel is finite dimensional, P has a finite rank and generates a smoothing operator. By spectral theory,

$$D^{-p} = \frac{1}{\Gamma(p/2)} \int_0^\infty t^{p/2} e^{-tD^2} t^{-1} dt = \frac{1}{\Gamma(p/2)} \left(\int_0^\epsilon + \int_\epsilon^\infty \right) t^{p/2} e^{-tD^2} t^{-1} dt$$

The second integral is a smooth operator since the map $x \rightarrow \int_\epsilon^\infty t^{p/2} e^{-tx^2} t^{-1} dt$ is in the Schwartz space \mathcal{S} .

Define the first integral as D_ϵ^{-p} and choose ϵ small enough such that for $0 < t \leq \epsilon$ and x and y close enough,

$$\left| k_t(x, y) - \frac{1}{(4\pi t)^{d/2}} \sum_{j=0}^{(d-p)/2} t^j \sqrt{\det g_x} k_j(x, y) e^{-d_g(x, y)^2/4t} \right| \leq c t^{p/2} e^{-d_g(x, y)^2/4t}.$$

Thus

$$\begin{aligned} \Gamma\left(\frac{d}{2}\right) \text{tr}(k_{D_\epsilon^{-p}}(x, y)) &= \int_0^\infty t^{p/2} \text{tr}(k_t(x, y)) t^{-1} dt \\ &= \frac{\sqrt{\det g_x}}{(4\pi)^{d/2}} \sum_{j=0}^{(d-p)/2} \text{tr}(k_j(x, y)) \int_0^\epsilon t^{j-(p-d)/2} e^{-d_g(x, y)^2/4t} dt \\ &\quad + \mathcal{O}\left(\int_0^\epsilon e^{-d_g(x, y)^2/4t} dt\right). \end{aligned}$$

For m integer and $\mu > 0$, we get after a change of variable $t \rightarrow t^{-1}$,

$$\int_0^\epsilon t^m e^{-\mu/t} t^{-1} dt = \mu^m \int_{\mu\epsilon}^\infty t^{-m} e^{-t} t^{-1} dt = \begin{cases} \text{Polynomial in } \frac{1}{\mu} + \mathcal{O}(1) & \text{for } m < 0, \\ -\log \mu + \mathcal{O}(1) & \text{for } m = 0, \\ \mathcal{O}(1) & \text{for } m > 0. \end{cases}$$

Thus, the logarithmic behavior of $\Gamma(\frac{d}{2}) \text{tr}(k_{D_\epsilon^{-p}}(x, y))$ comes from

$$\begin{aligned} \frac{\sqrt{\det g_x}}{(4\pi)^{d/2}} \text{tr}(k_{(d-p)/2}(x, y)) \int_0^\epsilon e^{-d_g(x,y)^2/4t} t^{-1} dt \\ = \frac{\sqrt{\det g_x}}{(4\pi)^{d/2}} \text{tr}(k_{(d-p)/2}(x, y)) \left(-\log(d_g(x, y)^2/4) + \mathcal{O}(1) \right) \\ = \frac{\sqrt{\det g_x}}{(4\pi)^{d/2}} \text{tr}(k_{(d-p)/2}(x, y)) \left(-2 \log(d_g(x, y)) + \mathcal{O}(1) \right). \end{aligned}$$

So

$$WRes(D^{-p}) = WRes(D_\epsilon^{-p}) = -\int_M c_{D_\epsilon^{-p}}(x) |dx| = \frac{2}{(4\pi)^{d/2}} \int_M \text{tr}(k_{(d-p)/2}(x, x)) \sqrt{\det g_x} |dx|,$$

which is, by definition, $\frac{2}{\Gamma(p/2)} a_{d-p}(D^2)$. \square

Few remarks are in order:

$$1) \text{ If } p = d, WRes(D^{-d}) = \frac{2}{\Gamma(p/2)} a_0(D^2) = \frac{2}{\Gamma(p/2)} \frac{\text{Rank}(E)}{(4\pi)^{d/2}} \text{Vol}(M).$$

Since $\text{Tr}(e^{-tD^2}) \underset{t \downarrow 0^+}{\sim} a_0(D^2) t^{-d/2}$, the Tauberian theorem used in Example 2.13 implies that

$D^{-d} = (D^{-2})^{d/2}$ is measurable and we obtain Connes' trace theorem 2.14

$$\text{Tr}_{Dix}(D^{-d}) = \text{Tr}_\omega(D^{-d}) = \frac{a_0(D^2)}{\Gamma(d/2+1)} = \frac{1}{d} WRes(D^{-d}).$$

2) When $D = \mathcal{D}$ and E is the spinor bundle, the Seeley-deWit coefficient $a_2(\mathcal{D}^2)$ (see (38) with $f = 1$) can be easily computed (see [44, 50]): if s is the scalar curvature,

$$a_2(\mathcal{D}^2) = -\frac{1}{12(4\pi)^{d/2}} \int_M s(x) d\text{vol}_g(x). \quad (40)$$

So $WRes(\mathcal{D}^{-d+2}) = \frac{2}{\Gamma(d/2-1)} a_2(\mathcal{D}^2) = c \int_M s(x) d\text{vol}_g(x)$. This is a quite important result since this last integral is nothing else but the Einstein–Hilbert action (70). In dimension 4, this is an example of invariant by diffeomorphisms, see (13).

5 Noncommutative integration

We already saw that the Wodzicki residue is a trace and, as such, can be viewed as an integral. But of course, it is quite natural to relate this integral to zeta functions used in (11): with notations of Section 1.4, let $P \in \Psi DO^{\mathbb{Z}}(M, E)$ and $D \in \Psi DO^1(M, E)$ which is elliptic. The definition of zeta function

$$\zeta_D^P(s) := \text{Tr} \left(P |D|^{-s} \right)$$

has been useful to prove that $WRes P = \text{Res}_{s=0} \zeta_D^P(s) = - \int_M c_P(x) |dx|$.

The aim now is to extend this notion to noncommutative spaces encoded in the notion of spectral triple.

References: [25, 31, 34, 37, 50].

5.1 Notion of spectral triple

The main properties of a compact spin Riemannian manifold M can be recaptured using the following triple $(\mathcal{A} = C^\infty(M), \mathcal{H} = L^2(M, S), \mathcal{D})$. The coordinates $x = (x^1, \dots, x^d)$ are exchanged with the algebra $C^\infty(M)$, the Dirac operator \mathcal{D} gives the dimension d as we saw in Theorem 3.9, but also the metric of M via Connes formula and more generally generates a quantized calculus. The idea of noncommutative geometry is to forget about the commutativity of the algebra and to impose axioms on a triplet $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ to generalize the above one in order to be able to obtain appropriate definitions of important notions: pseudodifferential operators, measure and integration theory, KO -theory, orientability, Poincaré duality, Hochschild (co)homology etc.

An important remark, probably due to Atiyah, is that the commutator of a pseudodifferential operator of order 1 (resp. order 0) with the multiplication by a function is a bounded operator (resp. compact). This is at the origin of the notion of Fredholm module (or K-cycle) with its K-homology class and via duality to its K-theory culminating with the Kasparov KK-theory. Thus, it is quite natural to define (unbounded) Fredholm module since for instance \mathcal{D} is unbounded:

Definition 5.1. *A spectral triple $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ is the data of an involutive (unital) algebra \mathcal{A} with a faithful representation π on a Hilbert space \mathcal{H} and a selfadjoint operator \mathcal{D} with compact resolvent (thus with discrete spectrum) such that $[\mathcal{D}, \pi(a)]$ is bounded for any $a \in \mathcal{A}$.*

We could impose the existence of a C^* -algebra A such that

$$\mathcal{A} := \{ a \in A \mid [\mathcal{D}, \pi(a)] \text{ is bounded} \}$$

is norm dense in A so \mathcal{A} is a pre- C^* -algebra stable by holomorphic calculus. Such \mathcal{A} is always a $*$ -subalgebra of A .

When there is no confusion, we will write a instead of $\pi(a)$.

We now give useful definitions:

Definition 5.2. *Let $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ be a spectral triple.*

It is even if there is a grading operator χ such that $\chi = \chi^$,*

$$[\chi, \pi(a)] = 0, \forall a \in \mathcal{A} \text{ and } \mathcal{D}\chi = -\chi\mathcal{D}.$$

It is real of KO -dimension $d \in \mathbb{Z}/8$ if there is an antilinear isometry $J : \mathcal{H} \rightarrow \mathcal{H}$ such that

$$J\mathcal{D} = \epsilon\mathcal{D}J, \quad J^2 = \epsilon', \quad J\chi = \epsilon''\chi J$$

with the following table for the signs $\epsilon, \epsilon', \epsilon''$

d	0	1	2	3	4	5	6	7
ϵ	1	-1	1	1	1	-1	1	1
ϵ'	1	1	-1	-1	-1	-1	1	1
ϵ''	1		-1		1		-1	

(41)

and the following commutation rules

$$[\pi(a), \pi(b)^\circ] = 0, \quad [[\mathcal{D}, \pi(a)], \pi(b)^\circ] = 0, \quad \forall a, b \in \mathcal{A} \quad (42)$$

where $\pi(a)^\circ := J\pi(a^*)J^{-1}$ is a representation of the opposite algebra \mathcal{A}° .

It is d -summable (or has metric dimension d) if the singular values of \mathcal{D} behave like $\mu_n(\mathcal{D}^{-1}) = \mathcal{O}(n^{-1/d})$.

It is regular if \mathcal{A} and $[\mathcal{D}, \mathcal{A}]$ are in the domain of δ^n for all $n \in \mathbb{N}$ where

$$\delta(T) := [|\mathcal{D}|, T].$$

It satisfies the finiteness condition if the space of smooth vectors $\mathcal{H}^\infty := \bigcap_k \text{Dom } \mathcal{D}^k$ is a finitely projective left \mathcal{A} -module.

It satisfies the orientation condition if there is a Hochschild cycle $c \in Z_d(\mathcal{A}, \mathcal{A} \otimes \mathcal{A}^\circ)$ such that $\pi_{\mathcal{D}}(c) = \chi$, where $\pi_{\mathcal{D}}((a \otimes b^\circ) \otimes a_1 \otimes \cdots \otimes a_d) := \pi(a)\pi(b)^\circ[\mathcal{D}, \pi(a_1)] \cdots [\mathcal{D}, \pi(a_d)]$ and d is its metric dimension.

The above definition of KO -dimension comes from the fact that a Dirac operator is a square root a Laplacian. This generates a sign problem which corresponds to a choice of a spin structure (or orientation). Up to some subtleties, the choice of a manifold of a chosen homotopy needs a Poincaré duality between homology and cohomology and the necessary refinement yields to the KO -homology introduced by Atiyah and Singer.

An interesting example of noncommutative space of non-zero KO -dimension is given by the finite part of the noncommutative standard model [21, 28, 31].

Moreover, the reality (or charge conjugation in the commutative case) operator J is related to the problem of the adjoint: If \mathcal{M} is a von Neumann algebra acting on the Hilbert space \mathcal{H} with a cyclic and separating vector $\xi \in \mathcal{H}$ (which means $\mathcal{M}\xi$ is dense in \mathcal{H} and $a\xi = 0$ implies $a = 0$, for $a \in \mathcal{M}$), then the closure S of the map: $a\xi \rightarrow a^*\xi$ has an unbounded extension to \mathcal{H} with a polar decomposition $S = J\Delta^{1/2}$ where $\Delta := S^*S$ is a positive operator and J is antilinear operator such that $J\mathcal{M}J^{-1} = \mathcal{M}'$, see Tomita theory in [103]. This explains the commutation relations (42). Moreover $\Delta^{it}\mathcal{M}\Delta^{-it} = \mathcal{M}$, a point related to Definition 5.5.

A fundamental point is that a reconstruction of the manifold is possible, starting only with a spectral triple where the algebra is commutative (see [29] for a more precise formulation, and also [94]):

Theorem 5.3. [29] *Given a commutative spectral triple $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ satisfying the above axioms, then there exists a compact spin^c manifold M such that $\mathcal{A} \simeq C^\infty(M)$ and \mathcal{D} is a Dirac operator.*

The manifold is known as a set, $M = \text{Sp}(\mathcal{A}) = \text{Sp}(A)$. Notice that \mathcal{D} is known only via its principal symbol, so is not unique. J encodes the nuance between spin and spin^c structures. The spectral action selects the Levi-Civita connection so *the* Dirac operator \mathcal{D} .

The way, the operator \mathcal{D} recaptures the original Riemannian metric g of M is via the Connes' distance:

Definition 5.4. *Given a spectral triple $(\mathcal{A}, \mathcal{H}, \mathcal{D})$,*

$$d(\phi_1, \phi_2) := \sup\{ |\phi_1(a) - \phi_2(a)| \mid \|[\mathcal{D}, \pi(a)]\| \leq 1, a \in \mathcal{A} \} \quad (43)$$

defines a distance (eventually infinite) between two states ϕ_1, ϕ_2 on the C^ -algebra A .*

In a commutative geometry, any point $x \in M$ defines a state via $\phi_x : a \in C^\infty(M) \rightarrow a(x) \in \mathbb{C}$. Since the geodesic distance is also given by

$$d_g(x, y) = \sup\{ |a(x) - a(y)| \mid a \in C^\infty(M), \|\text{grad } a\|_\infty \leq 1 \},$$

we get $d(x, y) = d_g(x, y)$ because $\|c(da)\| = \|\text{grad } a\|_\infty$. Recall that g is uniquely determined by its distance function by Myers–Steenrod theorem: if $\alpha : (M, g) \rightarrow (M', g')$ is a bijection such that $d_{g'}(\alpha(x), \alpha(y)) = d_g(x, y)$ for $x, y \in M$, then $g = \alpha^*g'$.

The role of \mathcal{D} is non only to provide a metric by (43), but its homotopy class represents the K -homology fundamental class of the noncommutative space \mathcal{A} .

It is known that one cannot hear the shape of a drum since the knowledge of the spectrum of a Laplacian does not determine the metric of the manifold, even if its conformal class is given [7]. But Theorem 5.3 shows that one can hear the shape of a spinorial drum (or better say, of a spectral triple) since the knowledge of the spectrum of the Dirac operator and the volume form, via its cohomological content, is sufficient to recapture the metric and spin structure. See however the more precise refinement made in [30]: for instance, if (M, g) is a compact oriented smooth Riemannian manifold, the spectral triple $(L^\infty(M), L^2(M, \wedge T^*M), \mathcal{D})$ where $\mathcal{D} = d + d^*$ is the signature operator (see example after definition 3.2) uniquely determines the manifold M .

5.2 Notion of pseudodifferential operators

Definition 5.5. *Let $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ be a spectral triple.*

For $t \in \mathbb{R}$ define the map $F_t : T \in \mathcal{B}(\mathcal{H}) \rightarrow e^{it|\mathcal{D}|} T e^{-it|\mathcal{D}|}$ and for $\alpha \in \mathbb{R}$

$$OP^0 := \{ T \mid t \rightarrow F_t(T) \in C^\infty(\mathbb{R}, \mathcal{B}(\mathcal{H})) \} \text{ is the set of operators of order } \leq 0,$$

$$OP^\alpha := \{ T \mid T|\mathcal{D}|^{-\alpha} \in OP^0 \} \text{ is the set of operators of order } \leq \alpha.$$

Moreover, we set

$$\delta(T) := [|\mathcal{D}|, T], \quad \nabla(T) := [\mathcal{D}^2, T].$$

For instance, $C^\infty(M) = OP^0 \cap L^\infty(M)$ and $L^\infty(M)$ is the von Neumann algebra generated by $\mathcal{A} = C^\infty(M)$.

Proposition 5.6. *Assume that $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ is regular so $\mathcal{A} \subset OP^0 = \bigcap_{k \geq 0} \text{Dom } \delta^k \subset \mathcal{B}(\mathcal{H})$. Then, for any $\alpha, \beta \in \mathbb{R}$,*

$$OP^\alpha OP^\beta \subset OP^{\alpha+\beta}, \quad OP^\alpha \subset OP^\beta \text{ if } \alpha \leq \beta, \quad \delta(OP^\alpha) \subset OP^\alpha, \quad \nabla(OP^\alpha) \subset OP^{\alpha+1}.$$

As an example, let us compute the order of $X = a|\mathcal{D}|[\mathcal{D}, b]\mathcal{D}^{-3}$: since the order of a is 0, of $|\mathcal{D}|$ is 1, of $[\mathcal{D}, b]$ is 0 and of \mathcal{D}^{-3} is -3, we get $X \in OP^{-2}$.

Definition 5.7. Let $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ be a spectral triple and $\mathcal{D}(\mathcal{A})$ be the polynomial algebra generated by \mathcal{A} , \mathcal{A}° , \mathcal{D} and $|\mathcal{D}|$.

Define the set of pseudodifferential operators as

$$\Psi(\mathcal{A}) := \{ T \mid \forall N \in \mathbb{N}, \exists P \in \mathcal{D}(\mathcal{A}), R \in OP^{-N}, p \in \mathbb{N} \text{ such that } T = P|\mathcal{D}|^{-p} + R \}$$

The idea behind this definition is that we want to work modulo the set $OP^{-\infty}$ of smoothing operators. This explains the presence of the arbitrary N and R . In the commutative case of a manifold M with spectral triple $(C^\infty(M), L^2(M, E), \mathcal{D})$ where $\mathcal{D} \in \text{Diff}^1(M, E)$, we get the natural inclusion $\Psi(C^\infty(M)) \subset \Psi DO(M, E)$.

The reader should be aware that Definition 5.7 is not exactly the same as in [31, 34, 50] since it pays attention to the reality operator J when it is present.

5.3 Zeta-functions and dimension spectrum

Definition 5.8. For $P \in \Psi^*(\mathcal{A})$, we define the zeta-function associated to P (and \mathcal{D}) by

$$\zeta_{\mathcal{D}}^P : s \in \mathbb{C} \rightarrow \text{Tr} \left(P|\mathcal{D}|^{-s} \right) \quad (44)$$

which makes sense since for $\Re(s) \gg 1$, $P|\mathcal{D}|^{-s} \in \mathcal{L}^1(\mathcal{H})$.

The dimension spectrum of $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ is the set $\{ \text{poles of } \zeta_{\mathcal{D}}^P(s) \mid P \in \Psi^*(\mathcal{A}) \}$. It is said simple if it contains poles of order at most one.

The noncommutative integral of P is defined by

$$\int P := \text{Res}_{s=0} \zeta_{\mathcal{D}}^P(s). \quad (45)$$

In (44), we assume \mathcal{D} invertible since otherwise, one can replace \mathcal{D} by the invertible operator $\mathcal{D} + P$, P being the projection on $\text{Ker } \mathcal{D}$. This change does not modify the computation of the integrals \int which follow since $\int X = 0$ when X is a trace-class operator.

The notion of dimension spectrum contains more informations than the usual dimension even for a manifold as we will see in Proposition 5.34.

Remark 5.9. If $Sp(\mathcal{A}, \mathcal{H}, \mathcal{D})$ denotes the set of all poles of the functions $s \mapsto \text{Tr} \left(P|\mathcal{D}|^{-s} \right)$ where P is any pseudodifferential operator, then, $Sd(\mathcal{A}, \mathcal{H}, \mathcal{D}) \subseteq Sp(\mathcal{A}, \mathcal{H}, \mathcal{D})$.

When $Sp(\mathcal{A}, \mathcal{H}, \mathcal{D}) = \mathbb{Z}$, $Sd(\mathcal{A}, \mathcal{H}, \mathcal{D}) = \{ n - k : k \in \mathbb{N}_0 \}$: indeed, if P is a pseudodifferential operator in OP^0 , and $q \in \mathbb{N}$ is such that $q > n$, $P|\mathcal{D}|^{-s}$ is in $OP^{-\Re(s)}$ so is trace-class for s in a neighborhood of q ; as a consequence, q cannot be a pole of $s \mapsto \text{Tr} \left(P|\mathcal{D}|^{-s} \right)$.

Due to the little difference of behavior between scalar and nonscalar pseudodifferential operators (i.e. when coefficients like $[\mathcal{D}, a]$, $a \in \mathcal{A}$ appears in P of Definition 5.7), it is convenient to also introduce

Definition 5.10. Let $\mathcal{D}_1(\mathcal{A})$ be the algebra generated by \mathcal{A} , $J\mathcal{A}J^{-1}$ and \mathcal{D} , and $\Psi_1(\mathcal{A})$ be the set of pseudodifferential operators constructed as before with $\mathcal{D}_1(\mathcal{A})$ instead of $\mathcal{D}(\mathcal{A})$. Note that $\Psi_1(\mathcal{A})$ is subalgebra of $\Psi(\mathcal{A})$.

Remark that $\Psi_1(\mathcal{A})$ does not necessarily contain operators such as $|\mathcal{D}|^k$ where $k \in \mathbb{Z}$ is odd. This algebra is similar to the one defined in [13].

5.4 One-forms and fluctuations of \mathcal{D}

The unitary group $\mathcal{U}(\mathcal{A})$ of \mathcal{A} gives rise to the automorphism $\alpha_u : a \in \mathcal{A} \rightarrow uau^* \in \mathcal{A}$. This defines the inner automorphisms group $\text{Inn}(\mathcal{A})$ which is a normal subgroup of the automorphisms $\text{Aut}(\mathcal{A}) := \{\alpha \in \text{Aut}(\mathcal{A}) \mid \alpha(\mathcal{A}) \subset \mathcal{A}\}$. For instance, in case of a gauge theory, the algebra $\mathcal{A} = C^\infty(M, M_n(\mathbb{C})) \simeq C^\infty(M) \otimes M_n(\mathbb{C})$ is typically used. Then, $\text{Inn}(\mathcal{A})$ is locally isomorphic to $\mathcal{G} = C^\infty(M, \text{PSU}(n))$. Since $\text{Aut}(C^\infty(M)) \simeq \text{Diff}(M)$, we get a complete parallel analogy between following two exact sequences:

$$\begin{array}{ccccccc} 1 & \longrightarrow & \text{Inn}(\mathcal{A}) & \longrightarrow & \text{Aut}(\mathcal{A}) & \longrightarrow & \text{Aut}(\mathcal{A})/\text{Inn}(\mathcal{A}) \longrightarrow 1, \\ 1 & \longrightarrow & \mathcal{G} & \longrightarrow & \mathcal{G} \times \text{Diff}(M) & \longrightarrow & \text{Diff}(M) \longrightarrow 1. \end{array}$$

This justifies that the internal symmetries of physics have to be replaced by the inner automorphisms.

We are looking for an equivalence relation between $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ and $(\mathcal{A}', \mathcal{H}', \mathcal{D}')$ giving rise to the same geometry. Of course, we could use unitary equivalence: there exists a unitary $U : \mathcal{H} \rightarrow \mathcal{H}'$ such that $\mathcal{D}' = U\mathcal{D}U^*$, $U\pi(a)U^* := \pi(\alpha(a))$ for some $\alpha \in \text{Aut}(\mathcal{A})$, and in the even real case $[U, \chi] = [U, J] = 0$. But this is not useful since it does not change the metric (43). So we need to vary not only \mathcal{D} but the algebra and its representation.

The appropriate framework for inner fluctuations of a spectral triple $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ is Morita equivalence that we describe now:

\mathcal{A} is Morita equivalent to \mathcal{B} if there is a finite projective right \mathcal{A} -module \mathcal{E} such that $\mathcal{B} \simeq \text{End}_{\mathcal{A}}(\mathcal{E})$. Thus \mathcal{B} acts on $\mathcal{H}' = \mathcal{E} \otimes_{\mathcal{A}} \mathcal{H}$ and \mathcal{H}' is endowed with scalar product $\langle r \otimes \eta, s \otimes \xi \rangle := \langle \eta, \pi(r|s)\xi \rangle$ where $(\cdot|\cdot)$ is a pairing $\mathcal{E} \times \mathcal{E} \rightarrow \mathcal{A}$ that is \mathcal{A} -linear in the second variable and satisfies $(r|s) = (s|r)^*$, $(r|sa) = (r|s)a$ and $(s|s) > 0$ for $0 \neq s \in \mathcal{E}$ (this can be seen as a \mathcal{A} -valued inner product).

A natural operator \mathcal{D}' associated to \mathcal{B} and \mathcal{H}' is a linear map $\mathcal{D}'(r \otimes \eta) = r \otimes \mathcal{D}\eta + (\nabla r)\eta$ where $\nabla : \mathcal{E} \rightarrow \mathcal{E} \otimes_{\mathcal{A}} \Omega_{\mathcal{D}}^1(\mathcal{A})$ is a linear map obeying to Leibniz rule $\nabla(ra) = (\nabla r)a + r \otimes [\mathcal{D}, a]$ for $r \in \mathcal{E}$, $a \in \mathcal{A}$ where we took the following

Definition 5.11. *Let $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ be a spectral triple. The set of one-forms is defined as*

$$\Omega_{\mathcal{D}}^1(\mathcal{A}) := \text{span}\{a db \mid a, b \in \mathcal{A}\}, \quad db := [\mathcal{D}, b].$$

It is a \mathcal{A} -bimodule.

Such ∇ is called a connection on \mathcal{E} and by a result of Cuntz–Quillen, only projective modules admit (universal) connections (see [50][Proposition 8.3]). Since we want \mathcal{D}' selfadjoint, ∇ must be hermitean with respect to \mathcal{D} which means: $\pi((r|\nabla s) - (\nabla r|s)) = [\mathcal{D}, \pi(r|s)]$.

In particular, when $\mathcal{E} = \mathcal{A}$ (any algebra is Morita equivalent to itself) and \mathcal{A} is regarded as a right \mathcal{A} -module, \mathcal{E} has a natural hermitean connection with respect to \mathcal{D} given by $Ad_{\mathcal{D}} : a \in \mathcal{A} \rightarrow [\mathcal{D}, a] \in \Omega_{\mathcal{D}}^1(\mathcal{A})$ and using the Leibniz rule, any another hermitean connection ∇ must verify: $\nabla a = Ad_{\mathcal{D}} a + Aa$ where $A = A^* \in \Omega_{\mathcal{D}}^1(\mathcal{A})$. So this process, which does not change neither the algebra \mathcal{A} nor the Hilbert space \mathcal{H} , gives a natural hermitean fluctuation of \mathcal{D} :

$$\mathcal{D} \rightarrow \mathcal{D}_A := \mathcal{D} + A \text{ with } A = A^* \in \Omega_{\mathcal{D}}^1(\mathcal{A}).$$

In conclusion, the Morita equivalent geometries for $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ keeping fixed \mathcal{A} and \mathcal{H} is an affine space modelled on the selfadjoint part of $\Omega_{\mathcal{D}}^1(\mathcal{A})$.

For instance, in commutative geometries, $\Omega_{\mathcal{D}}^1(C^\infty(M)) = \{c(da) \mid a \in C^\infty(M)\}$.

When a reality operator J exists, we also want $\mathcal{D}_A J = \epsilon J \mathcal{D}_A$, so we choose

$$\mathcal{D}_{\tilde{A}} := \mathcal{D} + \tilde{A}, \quad \tilde{A} := A + \epsilon J A J^{-1}, \quad A = A^*. \quad (46)$$

The next two results show that, with the same algebra \mathcal{A} and Hilbert space \mathcal{H} , a fluctuation of \mathcal{D} still give rise to a spectral triple $(\mathcal{A}, \mathcal{H}, \mathcal{D}_A)$ or $(\mathcal{A}, \mathcal{H}, \mathcal{D}_{\tilde{A}})$.

Lemma 5.12. *Let $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ be a spectral triple with a reality operator J and chirality χ . If $A \in \Omega_{\mathcal{D}}^1$ is a one-form, the fluctuated Dirac operator \mathcal{D}_A or $\mathcal{D}_{\tilde{A}}$ is an operator with compact resolvent, and in particular its kernel is a finite dimensional space. This space is invariant by J and χ .*

Proof. Let T be a bounded operator and let z be in the resolvent of $\mathcal{D} + T$ and z' be in the resolvent of \mathcal{D} . Then

$$(\mathcal{D} + T - z)^{-1} = (\mathcal{D} - z')^{-1} [1 - (T + z' - z)(\mathcal{D} + T - z)^{-1}].$$

Since $(\mathcal{D} - z')^{-1}$ is compact by hypothesis and since the term in bracket is bounded, $\mathcal{D} + T$ has a compact resolvent. Applying this to $T = A + \epsilon J A J^{-1}$, \mathcal{D}_A has a finite dimensional kernel (see for instance [66, Theorem 6.29]).

Since according to the dimension, $J^2 = \pm 1$, J commutes or anticommutes with χ , χ commutes with the elements in the algebra \mathcal{A} and $\mathcal{D}\chi = -\chi\mathcal{D}$, see (41), we get $\mathcal{D}_A\chi = -\chi\mathcal{D}_A$ and $\mathcal{D}_A J = \pm J \mathcal{D}_A$ which gives the result. \square

Note that $\mathcal{U}(\mathcal{A})$ acts on \mathcal{D} by $\mathcal{D} \rightarrow \mathcal{D}_u = u\mathcal{D}u^*$ leaving invariant the spectrum of \mathcal{D} . Since $\mathcal{D}_u = \mathcal{D} + u[\mathcal{D}, u^*]$ and in a C^* -algebra, any element a is a linear combination of at most four unitaries, Definition 5.11 is quite natural.

The inner automorphisms of a spectral triple correspond to inner fluctuation of the metric defined by (43).

One checks directly that a fluctuation of a fluctuation is a fluctuation and that the unitary group $\mathcal{U}(\mathcal{A})$ is gauge compatible for the adjoint representation:

Lemma 5.13. *Let $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ be a spectral triple (which is eventually real) and $A \in \Omega_{\mathcal{D}}^1(\mathcal{A})$, $A = A^*$.*

(i) *If $B \in \Omega_{\mathcal{D}_A}^1(\mathcal{A})$ (or $B \in \Omega_{\mathcal{D}_{\tilde{A}}}^1(\mathcal{A})$),*

$$\mathcal{D}_B = \mathcal{D}_C \text{ (or } \mathcal{D}_{\tilde{B}} = \mathcal{D}_{\tilde{C}}) \text{ with } C := A + B.$$

(ii) *Let $u \in \mathcal{U}(\mathcal{A})$. Then $U_u := uJuJ^{-1}$ is a unitary of \mathcal{H} such that*

$$U_u \mathcal{D}_{\tilde{A}} U_u^* = \mathcal{D}_{\widetilde{\gamma_u(A)}}, \text{ where } \gamma_u(A) := u[\mathcal{D}, u^*] + uAu^*.$$

Remark 5.14. *To be an inner fluctuation is not a symmetric relation. It can append that $\mathcal{D}_A = 0$ with $\mathcal{D} \neq 0$.*

Lemma 5.15. *Let $(\mathcal{A}, \mathcal{D}, \mathcal{H})$ be a spectral triple and $X \in \Psi(\mathcal{A})$. Then*

$$\int X^* = \overline{\int X}.$$

If the spectral triple is real, then, for $X \in \Psi(\mathcal{A})$, $JXJ^{-1} \in \Psi(\mathcal{A})$ and

$$\int JXJ^{-1} = \int X^* = \overline{\int X}.$$

Proof. The first result follows from (for $\Re s$ large enough, so the operators are traceable)

$$\mathrm{Tr}(X^*|\mathcal{D}|^{-s}) = \mathrm{Tr}\left(\left(|\mathcal{D}|^{-\bar{s}}X\right)^*\right) = \overline{\mathrm{Tr}\left(|\mathcal{D}|^{-\bar{s}}X\right)} = \overline{\mathrm{Tr}\left(X|\mathcal{D}|^{-\bar{s}}\right)}.$$

The second result is due to the anti-linearity of J , $\mathrm{Tr}(JYJ^{-1}) = \overline{\mathrm{Tr}(Y)}$, and $J|\mathcal{D}| = |\mathcal{D}|J$, so

$$\mathrm{Tr}(X|\mathcal{D}|^{-s}) = \overline{\mathrm{Tr}(JX|\mathcal{D}|^{-s}J^{-1})} = \overline{\mathrm{Tr}(JXJ^{-1}|\mathcal{D}|^{-\bar{s}})}. \quad \square$$

Corollary 5.16. *For any one-form $A = A^*$, and for $k, l \in \mathbb{N}$,*

$$\int A^l \mathcal{D}^{-k} \in \mathbb{R}, \quad \int (A\mathcal{D}^{-1})^k \in \mathbb{R}, \quad \int A^l |\mathcal{D}|^{-k} \in \mathbb{R}, \quad \int \chi A^l |\mathcal{D}|^{-k} \in \mathbb{R}, \quad \int A^l \mathcal{D} |\mathcal{D}|^{-k} \in \mathbb{R}.$$

We remark that the fluctuations leave invariant the first term of the spectral action (74). This is a generalization of the fact that in the commutative case, the noncommutative integral depends only on the principal symbol of the Dirac operator \mathcal{D} and this symbol is stable by adding a gauge potential like in $\mathcal{D} + A$. Note however that the symmetrized gauge potential $A + \epsilon JAJ^{-1}$ is always zero in this case for any selfadjoint one-form A , see (64).

Theorem 5.17. *Let $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ be a regular spectral triple which is simple and let $A \in \Omega_{\mathcal{D}}^1(\mathcal{A})$ be a selfadjoint gauge potential. Then,*

$$\zeta_{\mathcal{D}_{\tilde{A}}}(0) = \zeta_{\mathcal{D}}(0) + \sum_{q=1}^n \frac{(-1)^q}{q} \int (\tilde{A}\mathcal{D}^{-1})^q. \quad (47)$$

The proof needs few preliminaries.

Definition 5.18. *For an operator T , define the one-parameter group and notation*

$$\begin{aligned} \sigma_z(T) &:= |D|^z T |D|^{-z}, \quad z \in \mathbb{C}. \\ \epsilon(T) &:= \nabla(T) D^{-2}, \quad (\text{recall that } \nabla(T) = [\mathcal{D}^2, T]). \end{aligned}$$

The expansion of the one-parameter group σ_z gives for $T \in OP^q$

$$\sigma_z(T) \sim \sum_{r=0}^N g(z, r) \epsilon^r(T) \quad \text{mod } OP^{-N-1+q} \quad (48)$$

where $g(z, r) := \frac{1}{r!} \left(\frac{z}{2}\right) \cdots \left(\frac{z}{2} - (r-1)\right) = \binom{z/2}{r}$ with the convention $g(z, 0) := 1$.

We fix a regular spectral triple $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ of dimension d and a self-adjoint 1-form A . Despite previous remark before Lemma 5.15, we pay attention here to the kernel of \mathcal{D}_A since this operator can be non-invertible even if \mathcal{D} is, so we define

$$\begin{aligned} \mathcal{D}_A &:= \mathcal{D} + \tilde{A} \quad \text{where } \tilde{A} := A + \epsilon JAJ^{-1}, \\ D_A &:= \mathcal{D}_A + P_A \end{aligned} \quad (49)$$

where P_A is the projection on $\mathrm{Ker} \mathcal{D}_A$. Remark that $\tilde{A} \in \mathcal{D}(\mathcal{A}) \cap OP^0$ and $\mathcal{D}_A \in \mathcal{D}(\mathcal{A}) \cap OP^1$. We note

$$V_A := P_A - P_0.$$

As the following lemma shows, V_A is a smoothing operator:

Lemma 5.19. (i) $\bigcap_{k \geq 1} \text{Dom}(\mathcal{D}_A)^k \subseteq \bigcap_{k \geq 1} \text{Dom} |D|^k$.

(ii) $\text{Ker } \mathcal{D}_A \subseteq \bigcap_{k \geq 1} \text{Dom} |D|^k$.

(iii) For any $\alpha, \beta \in \mathbb{R}$, $|D|^\beta P_A |D|^\alpha$ is bounded.

(iv) $P_A \in OP^{-\infty}$.

Proof. (i) Let us define for any $p \in \mathbb{N}$, $R_p := (\mathcal{D}_A)^p - \mathcal{D}^p$, so $R_p \in OP^{p-1}$ and moreover $R_p(\text{Dom} |D|^p) \subseteq \text{Dom} |D|$.

Let us fix $k \in \mathbb{N}$, $k \geq 2$. Since $\text{Dom } \mathcal{D}_A = \text{Dom } \mathcal{D} = \text{Dom} |D|$, we have

$$\text{Dom}(\mathcal{D}_A)^k = \{ \phi \in \text{Dom} |D| : (\mathcal{D}^j + R_j) \phi \in \text{Dom} |D|, \forall j \ 1 \leq j \leq k-1 \}.$$

Let $\phi \in \text{Dom}(\mathcal{D}_A)^k$. We prove by recurrence that for any $j \in \{1, \dots, k-1\}$, $\phi \in \text{Dom} |D|^{j+1}$:

We have $\phi \in \text{Dom} |D|$ and $(\mathcal{D} + R_1) \phi \in \text{Dom} |D|$. Thus, since $R_1 \phi \in \text{Dom} |D|$, we have $\mathcal{D} \phi \in \text{Dom} |D|$, which proves that $\phi \in \text{Dom} |D|^2$. Hence, case $j = 1$ is done.

Suppose now $\phi \in \text{Dom} |D|^{j+1}$ for a $j \in \{1, \dots, k-2\}$. Since $(\mathcal{D}^{j+1} + R_{j+1}) \phi \in \text{Dom} |D|$, and $R_{j+1} \phi \in \text{Dom} |D|$, we get $\mathcal{D}^{j+1} \phi \in \text{Dom} |D|$, which proves that $\phi \in \text{Dom} |D|^{j+2}$.

Finally, if we set $j = k-1$, we get $\phi \in \text{Dom} |D|^k$, so $\text{Dom}(\mathcal{D}_A)^k \subseteq \text{Dom} |D|^k$.

(ii) follows from $\text{Ker } \mathcal{D}_A \subseteq \bigcap_{k \geq 1} \text{Dom}(\mathcal{D}_A)^k$ and (i).

(iii) Let us first check that $|D|^\alpha P_A$ is bounded. We define D_0 as the operator with domain $\text{Dom } D_0 = \text{Im } P_A \cap \text{Dom} |D|^\alpha$ and such that $D_0 \phi = |D|^\alpha \phi$. Since $\text{Dom } D_0$ is finite dimensional, D_0 extends as a bounded operator on \mathcal{H} with finite rank. We have

$$\sup_{\phi \in \text{Dom} |D|^\alpha P_A, \|\phi\| \leq 1} \||D|^\alpha P_A \phi\| \leq \sup_{\phi \in \text{Dom } D_0, \|\phi\| \leq 1} \||D|^\alpha \phi\| = \|D_0\| < \infty$$

so $|D|^\alpha P_A$ is bounded. We can remark that by (ii), $\text{Dom } D_0 = \text{Im } P_A$ and $\text{Dom} |D|^\alpha P_A = \mathcal{H}$.

Let us prove now that $P_A |D|^\alpha$ is bounded: Let $\phi \in \text{Dom } P_A |D|^\alpha = \text{Dom} |D|^\alpha$. By (ii), we have $\text{Im } P_A \subseteq \text{Dom} |D|^\alpha$ so we get

$$\begin{aligned} \|P_A |D|^\alpha \phi\| &\leq \sup_{\psi \in \text{Im } P_A, \|\psi\| \leq 1} | \langle \psi, |D|^\alpha \phi \rangle | \leq \sup_{\psi \in \text{Im } P_A, \|\psi\| \leq 1} | \langle |D|^\alpha \psi, \phi \rangle | \\ &\leq \sup_{\psi \in \text{Im } P_A, \|\psi\| \leq 1} \||D|^\alpha \psi\| \|\phi\| = \|D_0\| \|\phi\|. \end{aligned}$$

(iv) For any $k \in \mathbb{N}_0$ and $t \in \mathbb{R}$, $\delta^k(P_A) |D|^t$ is a linear combination of terms of the form $|D|^\beta P_A |D|^\alpha$, so the result follows from (iii). \square

Remark 5.20. We will see later on the noncommutative torus example how important is the difference between \mathcal{D}_A and $\mathcal{D} + A$. In particular, the inclusion $\text{Ker } \mathcal{D} \subseteq \text{Ker } \mathcal{D} + A$ is not satisfied since A does not preserve $\text{Ker } \mathcal{D}$ contrarily to \tilde{A} .

Let us define

$$\begin{aligned} X &:= \mathcal{D}_A^2 - \mathcal{D}^2 = \tilde{A} \mathcal{D} + \mathcal{D} \tilde{A} + \tilde{A}^2, \\ X_V &:= X + V_A, \end{aligned}$$

thus $X \in \mathcal{D}_1(\mathcal{A}) \cap OP^1$ and by Lemma 5.19,

$$X_V \sim X \pmod{OP^{-\infty}}. \tag{50}$$

We will use

$$Y := \log(D_A^2) - \log(D^2)$$

which makes sense since $D_A^2 = \mathcal{D}_A^2 + P_A$ is invertible for any A . By definition of X_V , we get

$$Y = \log(D^2 + X_V) - \log(D^2).$$

Lemma 5.21. (i) Y is a pseudodifferential operator in OP^{-1} with the following expansion for any $N \in \mathbb{N}$

$$Y \sim \sum_{p=1}^N \sum_{k_1, \dots, k_p=0}^{N-p} \frac{(-1)^{|k|_1+p+1}}{|k|_1+p} \nabla^{k_p} (X \nabla^{k_{p-1}} (\dots X \nabla^{k_1} (X) \dots)) D^{-2(|k|_1+p)} \quad \text{mod } OP^{-N-1}.$$

(ii) For any $N \in \mathbb{N}$ and $s \in \mathbb{C}$,

$$|D_A|^{-s} \sim |D|^{-s} + \sum_{p=1}^N K_p(Y, s) |D|^{-s} \quad \text{mod } OP^{-N-1-\Re(s)} \quad (51)$$

with $K_p(Y, s) \in OP^{-p}$.

Proof. (i) We follow [13, Lemma 2.2]. By functional calculus, $Y = \int_0^\infty I(\lambda) d\lambda$, where

$$I(\lambda) \sim \sum_{p=1}^N (-1)^{p+1} \left((D^2 + \lambda)^{-1} X_V \right)^p (D^2 + \lambda)^{-1} \quad \text{mod } OP^{-N-3}.$$

By (50), $\left((D^2 + \lambda)^{-1} X_V \right)^p \sim \left((D^2 + \lambda)^{-1} X \right)^p \quad \text{mod } OP^{-\infty}$ and we get

$$I(\lambda) \sim \sum_{p=1}^N (-1)^{p+1} \left((D^2 + \lambda)^{-1} X \right)^p (D^2 + \lambda)^{-1} \quad \text{mod } OP^{-N-3}.$$

We set $A_p(X) := \left((D^2 + \lambda)^{-1} X \right)^p (D^2 + \lambda)^{-1}$ and $L := (D^2 + \lambda)^{-1} \in OP^{-2}$ for a fixed λ . Since $[D^2 + \lambda, X] \sim \nabla(X) \quad \text{mod } OP^{-\infty}$, a recurrence proves that if T is an operator in OP^r , then, for $q \in \mathbb{N}_0$,

$$A_1(T) = LTL \sim \sum_{k=0}^q (-1)^k \nabla^k(T) L^{k+2} \quad \text{mod } OP^{r-q-5}.$$

With $A_p(X) = LX A_{p-1}(X)$, another recurrence gives, for any $q \in \mathbb{N}_0$,

$$A_p(X) \sim \sum_{k_1, \dots, k_p=0}^q (-1)^{|k|_1} \nabla^{k_p} (X \nabla^{k_{p-1}} (\dots X \nabla^{k_1} (X) \dots)) L^{|k|_1+p+1} \quad \text{mod } OP^{-q-p-3},$$

which entails that

$$I(\lambda) \sim \sum_{p=1}^N (-1)^{p+1} \sum_{k_1, \dots, k_p=0}^{N-p} (-1)^{|k|_1} \nabla^{k_p} (X \nabla^{k_{p-1}} (\dots X \nabla^{k_1} (X) \dots)) L^{|k|_1+p+1} \quad \text{mod } OP^{-N-3}.$$

With $\int_0^\infty (D^2 + \lambda)^{-(|k|_1+p+1)} d\lambda = \frac{1}{|k|_1+p} D^{-2(|k|_1+p)}$, we get the result provided we control the remainders. Such a control is given in [13, (2.27)].

(ii) Applied to $|D_A|^{-s} = e^{B-(s/2)Y} e^{-B} |D|^{-s}$ where $B := (-s/2) \log(D^2)$, the Duhamel's expansion formula

$$e^{U+V} e^{-U} = \sum_{n=0}^{\infty} \int_{0 \leq t_1 \leq t_2 \leq \dots \leq t_n \leq 1} V(t_1) \cdots V(t_n) dt_1 \cdots dt_n$$

with $V(t) := e^{tU} V e^{-tU}$ gives

$$|D_A|^{-s} = |D|^{-s} + \sum_{p=1}^{\infty} K_p(Y, s) |D|^{-s}. \quad (52)$$

and each $K_p(Y, s)$ is in OP^{-p} . □

Corollary 5.22. *For any $p \in \mathbb{N}$ and $r_1, \dots, r_p \in \mathbb{N}_0$, $\varepsilon^{r_1}(Y) \cdots \varepsilon^{r_p}(Y) \in \Psi_1(\mathcal{A})$.*

Proof. If for any $q \in \mathbb{N}$ and $k = (k_1, \dots, k_q) \in \mathbb{N}_0^q$,

$$\Gamma_q^k(X) := \frac{(-1)^{|k|_1+q+1}}{|k|_1+q} \nabla^{k_q} (X \nabla^{k_{q-1}} (\cdots X \nabla^{k_1} (X) \cdots)),$$

then, $\Gamma_q^k(X) \in OP^{|k|_1+q}$. For any $N \in \mathbb{N}$,

$$Y \sim \sum_{q=1}^N \sum_{k_1, \dots, k_q=0}^{N-q} \Gamma_q^k(X) D^{-2(|k|_1+q)} \quad \text{mod } OP^{-N-1}. \quad (53)$$

Since the $\Gamma_q^k(X)$ are in $\mathcal{D}(\mathcal{A})$, this proves with (53) that Y and thus $\varepsilon^r(Y) = \nabla^r(Y) D^{-2r}$, are also in $\Psi_1(\mathcal{A})$. □

Proof of Theorem 5.17. Again, we follow [13]. Since the spectral triple is simple, equation (52) entails that

$$\zeta_{D_A}(0) - \zeta_D(0) = \text{Tr}(K_1(Y, s) |D|^{-s})|_{s=0}.$$

Thus, with (48), we get $\zeta_{D_A}(0) - \zeta_D(0) = -\frac{1}{2} f Y$.

Now the conclusion follows from $f \log((1+S)(1+T)) = f \log(1+S) + f \log(1+T)$ for $S, T \in \Psi(\mathcal{A}) \cap OP^{-1}$ (since $\log(1+S) = \sum_{n=1}^{\infty} \frac{(-1)^{n+1}}{n} S^n$) with $S = D^{-1}A$ and $T = AD^{-1}$; so $f \log(1+XD^{-2}) = 2f \log(1+AD^{-1})$ and

$$-\frac{1}{2} f Y = \sum_{q=1}^n \frac{(-1)^q}{q} f (\tilde{A} D^{-1})^q. \quad \square$$

Lemma 5.23. *For any $k \in \mathbb{N}_0$,*

$$\text{Res}_{s=d-k} \zeta_{D_A}(s) = \text{Res}_{s=d-k} \zeta_D(s) + \sum_{p=1}^k \sum_{r_1, \dots, r_p=0}^{k-p} \text{Res}_{s=d-k} h(s, r, p) \text{Tr}(\varepsilon^{r_1}(Y) \cdots \varepsilon^{r_p}(Y) |D|^{-s}),$$

where

$$h(s, r, p) := (-s/2)^p \int_{0 \leq t_1 \leq \dots \leq t_p \leq 1} g(-st_1, r_1) \cdots g(-st_p, r_p) dt.$$

Proof. By Lemma 5.21 (ii), $|D_A|^{-s} \sim |D|^{-s} + \sum_{p=1}^k K_p(Y, s)|D|^{-s} \pmod{OP^{-(k+1)-\Re(s)}}$, where the convention $\sum_{\emptyset} = 0$ is used. Thus, we get for s in a neighborhood of $d - k$,

$$|D_A|^{-s} - |D|^{-s} - \sum_{p=1}^k K_p(Y, s)|D|^{-s} \in OP^{-(k+1)-\Re(s)} \subseteq \mathcal{L}^1(\mathcal{H})$$

which gives

$$\operatorname{Res}_{s=d-k} \zeta_{D_A}(s) = \operatorname{Res}_{s=d-k} \zeta_D(s) + \sum_{p=1}^k \operatorname{Res}_{s=d-k} \operatorname{Tr} \left(K_p(Y, s)|D|^{-s} \right). \quad (54)$$

Let us fix $1 \leq p \leq k$ and $N \in \mathbb{N}$. By (48) we get

$$K_p(Y, s) \sim \left(-\frac{s}{2}\right)^p \int_{0 \leq t_1 \leq \dots \leq t_p \leq 1} \sum_{r_1, \dots, r_p=0}^N g(-st_1, r_1) \cdots g(-st_p, r_p) \varepsilon^{r_1}(Y) \cdots \varepsilon^{r_p}(Y) dt \pmod{OP^{-N-p-1}}. \quad (55)$$

If we now take $N = k - p$, we get for s in a neighborhood of $d - k$

$$K_p(Y, s)|D|^{-s} - \sum_{r_1, \dots, r_p=0}^{k-p} h(s, r, p) \varepsilon^{r_1}(Y) \cdots \varepsilon^{r_p}(Y)|D|^{-s} \in OP^{-k-1-\Re(s)} \subseteq \mathcal{L}^1(\mathcal{H})$$

so (54) gives the result. \square

Our operators $|D_A|^k$ are pseudodifferential operators:

Lemma 5.24. *For any $k \in \mathbb{Z}$, $|D_A|^k \in \Psi^k(\mathcal{A})$.*

Proof. Using (55), we see that $K_p(Y, s)$ is a pseudodifferential operator in OP^{-p} , so (51) proves that $|D_A|^k$ is a pseudodifferential operator in OP^k . \square

The following result is quite important since it shows that one can use f for D or D_A :

Proposition 5.25. *If the spectral triple is simple, $\operatorname{Res}_{s=0} \operatorname{Tr} \left(P|D_A|^{-s} \right) = fP$ for any pseudodifferential operator P . In particular, for any $k \in \mathbb{N}_0$*

$$\int |D_A|^{-(d-k)} = \operatorname{Res}_{s=d-k} \zeta_{D_A}(s).$$

Proof. Suppose $P \in OP^k$ with $k \in \mathbb{Z}$ and let us fix $p \geq 1$. With (55), we see that for any $N \in \mathbb{N}$,

$$PK_p(Y, s)|D|^{-s} \sim \sum_{r_1, \dots, r_p=0}^N h(s, r, p) P \varepsilon^{r_1}(Y) \cdots \varepsilon^{r_p}(Y)|D|^{-s} \pmod{OP^{-N-p-1+k-\Re(s)}}.$$

Thus if we take $N = d - p + k$, we get

$$\operatorname{Res}_{s=0} \operatorname{Tr} \left(PK_p(Y, s)|D|^{-s} \right) = \sum_{r_1, \dots, r_p=0}^{n-p+k} \operatorname{Res}_{s=0} h(s, r, p) \operatorname{Tr} \left(P \varepsilon^{r_1}(Y) \cdots \varepsilon^{r_p}(Y)|D|^{-s} \right).$$

Since $s = 0$ is a zero of the analytic function $s \mapsto h(s, r, p)$ and $s \mapsto \text{Tr } P\varepsilon^{r_1}(Y) \cdots \varepsilon^{r_p}(Y)|D|^{-s}$ has only simple poles by hypothesis, we get $\text{Res}_{s=0} h(s, r, p) \text{Tr} \left(P\varepsilon^{r_1}(Y) \cdots \varepsilon^{r_p}(Y)|D|^{-s} \right) = 0$ and

$$\text{Res}_{s=0} \text{Tr} \left(PK_p(Y, s)|D|^{-s} \right) = 0. \quad (56)$$

Using (51), $P|D_A|^{-s} \sim P|D|^{-s} + \sum_{p=1}^{k+d} PK_p(Y, s)|D|^{-s} \pmod{OP^{-d-1-\Re(s)}}$ and thus,

$$\text{Res}_{s=0} \text{Tr}(P|D_A|^{-s}) = \int P + \sum_{p=1}^{k+d} \text{Res}_{s=0} \text{Tr} \left(PK_p(Y, s)|D|^{-s} \right). \quad (57)$$

The result now follows from (56) and (57). To get the last equality, one uses the pseudodifferential operator $|D_A|^{-(d-k)}$. \square

Proposition 5.26. *If the spectral triple is simple, then*

$$\int |D_A|^{-d} = \int |D|^{-d}. \quad (58)$$

Proof. Lemma 5.23 and previous proposition for $k = 0$. \square

Lemma 5.27. *If the spectral triple is simple,*

- (i) $\int |D_A|^{-(d-1)} = \int |D|^{-(d-1)} - \left(\frac{d-1}{2}\right) \int X|D|^{-d-1}$.
- (ii) $\int |D_A|^{-(d-2)} = \int |D|^{-(d-2)} + \frac{d-2}{2} \left(-\int X|D|^{-d} + \frac{d}{4} \int X^2|D|^{-2-d} \right)$.

Proof. (i) By (51),

$$\text{Res}_{s=d-1} \zeta_{D_A}(s) - \zeta_D(s) = \text{Res}_{s=d-1} (-s/2) \text{Tr} \left(Y|D|^{-s} \right) = -\frac{d-1}{2} \text{Res}_{s=0} \text{Tr} \left(Y|D|^{-(d-1)}|D|^{-s} \right)$$

where for the last equality we use the simple dimension spectrum hypothesis. Lemma 5.21 (i) yields $Y \sim XD^{-2} \pmod{OP^{-2}}$ and $Y|D|^{-(d-1)} \sim X|D|^{-d-1} \pmod{OP^{-d-1}} \subseteq \mathcal{L}^1(\mathcal{H})$. Thus,

$$\text{Res}_{s=0} \text{Tr} \left(Y|D|^{-(d-1)}|D|^{-s} \right) = \text{Res}_{s=0} \text{Tr} \left(X|D|^{-d-1}|D|^{-s} \right) = \int X|D|^{-d-1}.$$

(ii) Lemma 5.23 (ii) gives

$$\text{Res}_{s=d-2} \zeta_{D_A}(s) = \text{Res}_{s=d-2} \zeta_D(s) + \text{Res}_{s=d-2} \sum_{r=0}^1 h(s, r, 1) \text{Tr} \left(\varepsilon^r(Y)|D|^{-s} \right) + h(s, 0, 2) \text{Tr} \left(Y^2|D|^{-s} \right).$$

We have $h(s, 0, 1) = -\frac{s}{2}$, $h(s, 1, 1) = \frac{1}{2}(\frac{s}{2})^2$ and $h(s, 0, 2) = \frac{1}{2}(\frac{s}{2})^2$. Using again Lemma 5.21 (i),

$$Y \sim XD^{-2} - \frac{1}{2}\nabla(X)D^{-4} - \frac{1}{2}X^2D^{-4} \pmod{OP^{-3}}.$$

Thus,

$$\text{Res}_{s=d-2} \text{Tr} \left(Y|D|^{-s} \right) = \int X|D|^{-d} - \frac{1}{2} \int (\nabla(X) + X^2)|D|^{-2-d}.$$

Moreover, using $\int \nabla(X)|D|^{-k} = 0$ for any $k \geq 0$ since \int is a trace,

$$\text{Res}_{s=d-2} \text{Tr} \left(\varepsilon(Y)|D|^{-s} \right) = \text{Res}_{s=d-2} \text{Tr} \left(\nabla(X)D^{-4}|D|^{-s} \right) = \int \nabla(X)|D|^{-2-d} = 0.$$

Similarly, since $Y \sim XD^{-2} \pmod{OP^{-2}}$ and $Y^2 \sim X^2D^{-4} \pmod{OP^{-3}}$, we get

$$\operatorname{Res}_{s=d-2} \operatorname{Tr} \left(Y^2 |D|^{-s} \right) = \operatorname{Res}_{s=d-2} \operatorname{Tr} \left(X^2 D^{-4} |D|^{-s} \right) = \int X^2 |D|^{-2-d}.$$

Thus,

$$\begin{aligned} \operatorname{Res}_{s=d-2} \zeta_{D_A}(s) &= \operatorname{Res}_{s=d-2} \zeta_D(s) + \left(-\frac{d-2}{2}\right) \left(\int X |D|^{-d} - \frac{1}{2} \int (\nabla(X) + X^2) |D|^{-2-d} \right) \\ &\quad + \frac{1}{2} \left(\frac{d-2}{2}\right)^2 \int \nabla(X) |D|^{-2-d} + \frac{1}{2} \left(\frac{d-2}{2}\right)^2 \int X^2 |D|^{-2-d}. \end{aligned}$$

Finally,

$$\operatorname{Res}_{s=d-2} \zeta_{D_A}(s) = \operatorname{Res}_{s=d-2} \zeta_D(s) + \left(-\frac{d-2}{2}\right) \left(\int X |D|^{-d} - \frac{1}{2} \int X^2 |D|^{-2-d} \right) + \frac{1}{2} \left(\frac{d-2}{2}\right)^2 \int X^2 |D|^{-2-d}$$

and the result follows from Proposition 5.25. \square

Corollary 5.28. *If the spectral triple satisfies $\int |D|^{-(d-2)} = \int \tilde{A} \mathcal{D} |D|^{-d} = \int \mathcal{D} \tilde{A} |D|^{-d} = 0$, then*

$$\int |D_A|^{-(d-2)} = \frac{d(d-2)}{4} \left(\int \tilde{A} \mathcal{D} \tilde{A} \mathcal{D} |D|^{-d-2} + \frac{d-2}{d} \int \tilde{A}^2 |D|^{-d} \right).$$

Proof. By previous lemma,

$$\operatorname{Res}_{s=d-2} \zeta_{D_A}(s) = \frac{d-2}{2} \left(- \int \tilde{A}^2 |D|^{-d} + \frac{d}{4} \int (\tilde{A} \mathcal{D} \tilde{A} \mathcal{D} + \mathcal{D} \tilde{A} \mathcal{D} \tilde{A} + \tilde{A} \mathcal{D}^2 \tilde{A} + \mathcal{D} \tilde{A}^2 \mathcal{D}) |D|^{-d-2} \right).$$

Since $\nabla(\tilde{A}) \in OP^1$, the trace property of \int yields the result. \square

5.5 Tadpole

In [31], the following definition is introduced:

Definition 5.29. *In $(\mathcal{A}, \mathcal{H}, \mathcal{D})$, the tadpole $\operatorname{Tad}_{\mathcal{D}+A}(k)$ of order k , for $k \in \{d-l : l \in \mathbb{N}\}$ is the term linear in $A = A^* \in \Omega_{\mathcal{D}}^1$, in the Λ^k term of (74) (considered as an infinite series) where $\mathcal{D} \rightarrow \mathcal{D} + A$.*

If moreover, the triple $(\mathcal{A}, \mathcal{H}, \mathcal{D}, J)$ is real, the tadpole $\operatorname{Tad}_{\mathcal{D}+\tilde{A}}(k)$ is the term linear in A , in the Λ^k term of (74) where $\mathcal{D} \rightarrow \mathcal{D} + \tilde{A}$.

Proposition 5.30. *Let $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ be a spectral triple of dimension d with simple dimension spectrum. Then*

$$\operatorname{Tad}_{\mathcal{D}+A}(d-k) = -(d-k) \int A \mathcal{D} |D|^{-(d-k)-2}, \quad \forall k \neq d, \quad (59)$$

$$\operatorname{Tad}_{\mathcal{D}+A}(0) = - \int A \mathcal{D}^{-1}. \quad (60)$$

Moreover, if the triple is real, $\operatorname{Tad}_{\mathcal{D}+\tilde{A}} = 2 \operatorname{Tad}_{\mathcal{D}+A}$.

Proof. We already proved the following formula, for any $k \in \mathbb{N}$,

$$\int |\mathcal{D}_A|^{-(d-k)} = \int |\mathcal{D}|^{-(d-k)} + \sum_{p=1}^k \sum_{r_1, \dots, r_p=0}^{k-p} \text{Res}_{s=d-k} h(s, r, p) \text{Tr} \left(\varepsilon^{r_1}(Y) \cdots \varepsilon^{r_p}(Y) |\mathcal{D}|^{-s} \right),$$

with here $X := \tilde{A}\mathcal{D} + \mathcal{D}\tilde{A} + \tilde{A}^2$, $\tilde{A} := A + \epsilon JAJ^{-1}$.

As a consequence, for $k \neq n$, only the terms with $p = 1$ contribute to the linear part:

$$\text{Tad}_{\mathcal{D}+\tilde{A}}(d-k) = \text{Lin}_A \left(\int |\mathcal{D}_A|^{-(d-k)} \right) = \sum_{r=0}^{k-1} \text{Res}_{s=d-k} h(s, r, 1) \text{Tr} \left(\varepsilon^r(\text{Lin}_A(Y)) |\mathcal{D}|^{-s} \right).$$

We check that for any $N \in \mathbb{N}^*$,

$$\text{Lin}_A(Y) \sim \sum_{l=0}^{N-1} \Gamma_1^l(\tilde{A}\mathcal{D} + \mathcal{D}\tilde{A}) \mathcal{D}^{-2(l+1)} \quad \text{mod } OP^{-N-1}.$$

Since $\Gamma_1^l(\tilde{A}\mathcal{D} + \mathcal{D}\tilde{A}) = \frac{(-1)^l}{l+1} \nabla^l(\tilde{A}\mathcal{D} + \mathcal{D}\tilde{A}) = \frac{(-1)^l}{l+1} \{\nabla^l(\tilde{A}), \mathcal{D}\}$, we get, assuming the dimension spectrum to be simple

$$\begin{aligned} \text{Tad}_{\mathcal{D}+\tilde{A}}(d-k) &= \sum_{r=0}^{k-1} \text{Res}_{s=d-k} h(s, r, p) \text{Tr} \left(\varepsilon^r(\text{Lin}_A(Y)) |\mathcal{D}|^{-s} \right) \\ &= \sum_{r=0}^{k-1} h(n-k, r, 1) \sum_{l=0}^{k-1-r} \frac{(-1)^l}{l+1} \text{Res}_{s=d-k} \text{Tr} \left(\varepsilon^r(\{\nabla^l(\tilde{A}), \mathcal{D}\}) |\mathcal{D}|^{-s-2(l+1)} \right) \\ &= 2 \sum_{r=0}^{k-1} h(d-k, r, 1) \sum_{l=0}^{k-1-r} \frac{(-1)^l}{l+1} \int \nabla^{r+l}(\tilde{A}) \mathcal{D} |\mathcal{D}|^{-(d-k+2(r+l))-2} \\ &= -(n-k) \int \tilde{A} \mathcal{D} |\mathcal{D}|^{-(d-k)-2}, \end{aligned}$$

because in the last sum it remains only the case $r+l=0$, so $r=l=0$.

Formula (60) is a direct application of Theorem 5.17.

The link between $\text{Tad}_{\mathcal{D}+\tilde{A}}$ and $\text{Tad}_{\mathcal{D}+A}$ follows from $J\mathcal{D} = \epsilon \mathcal{D}J$ and Lemma 5.15. \square

Corollary 5.31. *In a real spectral triple $(\mathcal{A}, \mathcal{H}, \mathcal{D})$, if $A = A^* \in \Omega_{\mathcal{D}}^1(\mathcal{A})$ is such that $\tilde{A} = 0$, then $\text{Tad}_{\mathcal{D}+A}(k) = 0$ for any $k \in \mathbb{Z}$, $k \leq d$.*

The vanishing tadpole of order 0 has the following equivalence (see [13])

$$\int A \mathcal{D}^{-1} = 0, \forall A \in \Omega_{\mathcal{D}}^1(\mathcal{A}) \iff \int ab = \int a\alpha(b), \forall a, b \in \mathcal{A}, \quad (61)$$

where $\alpha(b) := \mathcal{D}b\mathcal{D}^{-1}$.

The existence of tadpoles is important since, for instance, $A = 0$ is not necessarily a stable solution of the classical field equation deduced from spectral action expansion, [51].

5.6 Commutative geometry

Definition 5.32. Consider a commutative spectral triple given by a compact Riemannian spin manifold M of dimension d without boundary and its Dirac operator \mathcal{D} associated to the Levi-Civita connection. This means $(\mathcal{A} := C^\infty(M), \mathcal{H} := L^2(M, S), \mathcal{D})$ where S is the spinor bundle over M . This triple is real since, due to the existence of a spin structure, the charge conjugation operator generates an anti-linear isometry J on \mathcal{H} such that

$$J a J^{-1} = a^*, \quad \forall a \in \mathcal{A}, \quad (62)$$

and when d is even, the grading is given by the chirality matrix

$$\chi := (-i)^{d/2} \gamma^1 \gamma^2 \cdots \gamma^d. \quad (63)$$

Such triple is said to be a commutative geometry.

In the polynomial algebra $\mathcal{D}(\mathcal{A})$ of Definition 5.7, we added \mathcal{A}° . In the commutative case, $\mathcal{A}^\circ \simeq J \mathcal{A} J^{-1} \simeq \mathcal{A}$ as indicated by (62) which also gives

$$J A J^{-1} = -\epsilon A^*, \quad \forall A \in \Omega_{\mathcal{D}}^1(\mathcal{A}) \text{ or } \tilde{A} = 0 \text{ when } A = A^*. \quad (64)$$

As noticed by Wodzicki, $\int P$ is equal to -2 times the coefficient in $\log t$ of the asymptotics of $\text{Tr}(P e^{-t \mathcal{D}^2})$ as $t \rightarrow 0$. It is remarkable that this coefficient is independent of \mathcal{D} as seen in Theorem 1.22 and this gives a close relation between the ζ function and heat kernel expansion with $WRes$. Actually, by [48, Theorem 2.7]

$$\text{Tr}(P e^{-t \mathcal{D}^2}) \underset{t \downarrow 0^+}{\sim} \sum_{k=0}^{\infty} a_k t^{(k - \text{ord}(P) - d)/2} + \sum_{k=0}^{\infty} (-a'_k \log t + b_k) t^k, \quad (65)$$

so

$$\int P = 2a'_0.$$

Remark that $\int, WRes$ are traces on $\Psi(C^\infty(M))$, thus Corollary 1.23 implies

$$\int P = c WRes P \quad (66)$$

Since, via Mellin transform, $\text{Tr}(P \mathcal{D}^{-2s}) = \frac{1}{\Gamma(s)} \int_0^\infty t^{s-1} \text{Tr}(P e^{-t \mathcal{D}^2}) dt$, the non-zero coefficient a'_k , $k \neq 0$ creates a pole of $\text{Tr}(P \mathcal{D}^{-2s})$ of order $k + 2$ because we get $\int_0^1 t^{s-1} \log(t)^k dt = \frac{(-1)^k k!}{s^{k+1}}$ and

$$\Gamma(s) = \frac{1}{s} + \gamma + s g(s) \quad (67)$$

where γ is the Euler constant and the function g is also holomorphic around zero.

We have $\int 1 = 0$ and more generally, $WRes(P) = 0$ for all zero-order pseudodifferential projections [112].

As the following remark shows, being a commutative geometry is more than just having a commutative algebra:

Remark 5.33. Since $J\pi(a)J^{-1} = \pi(a^*)$ for all $a \in \mathcal{A}$ and $\tilde{A} = 0$ for all $A = A^* \in \Omega_{\mathcal{D}}^1$ when \mathcal{A} is commutative by (64), one can only use $\mathcal{D}_A = \mathcal{D} + A$ to get fluctuation of \mathcal{D} : It is amazing to see that in the context of noncommutative geometry, to get an abelian gauge field, we need to go outside of abelian algebras. In particular, as pointed out in [107], a commutative manifold could support relativity but not electromagnetism.

However, we can have \mathcal{A} commutative and $J\pi(a)J^{-1} \neq \pi(a^*)$ for some $a \in \mathcal{A}$ [27, 68]: Let $\mathcal{A}_1 = \mathbb{C} \oplus \mathbb{C}$ represented on $\mathcal{H}_1 = \mathbb{C}^3$ with, for some complex number $m \neq 0$,

$$\pi_1(a) := \begin{pmatrix} b_1 & 0 & 0 \\ 0 & b_1 & 0 \\ 0 & 0 & b_2 \end{pmatrix}, \text{ for } a = (b_1, b_2) \in \mathcal{A},$$

and

$$\mathcal{D}_1 := \begin{pmatrix} 0 & m & m \\ \bar{m} & 0 & 0 \\ \bar{m} & 0 & 0 \end{pmatrix}, \chi_1 := \begin{pmatrix} 1 & 0 & 0 \\ 0 & -1 & 0 \\ 0 & 0 & -1 \end{pmatrix}, J_1 := \begin{pmatrix} 1 & 0 & 0 \\ 0 & 0 & 1 \\ 0 & 1 & 0 \end{pmatrix} \circ cc$$

where cc is the complex conjugation. Then $(\mathcal{A}_1, \mathcal{H}_1, \mathcal{D}_1)$ is a commutative real spectral triple of dimension $d = 0$ with non zero one-forms and such that $J_1\pi_1(a)J_1^{-1} = \pi_1(a^*)$ only if $a = (b_1, b_1)$.

Take now a commutative geometry $(\mathcal{A}_2 = C^\infty(M), \mathcal{H} = L^2(M, S), \mathcal{D}_2, \chi_2, J_2)$ defined in 5.32 where $d = \dim M$ is even, and then take the tensor product of the two spectral triples, namely $\mathcal{A} = \mathcal{A}_1 \otimes \mathcal{A}_2$, $\mathcal{H} = \mathcal{H}_1 \otimes \mathcal{H}_2$, $\pi = \pi_1 \otimes \pi_2$, $\mathcal{D} = \mathcal{D}_1 \otimes \chi_2 + 1 \otimes \mathcal{D}_2$, $\chi = \chi_1 \otimes \chi_2$ and J is either $\chi_1 J_1 \otimes J_2$ when $d \in \{2, 6\} \bmod 8$ or $J_1 \otimes J_2$ in the other cases, see [27, 105].

Then $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ is a real commutative spectral triple of dimension d such that $\tilde{A} \neq 0$ for some selfadjoint one-forms A , so is not exactly like in Definition 5.32.

Proposition 5.34. Let $Sp(M)$ be the dimension spectrum of a commutative geometry of dimension d . Then $Sp(M)$ is simple and $Sp(M) = \{d - k \mid k \in \mathbb{N}\}$.

Proof. Let $a \in \mathcal{A} = C^\infty(M)$ such that its trace norm $\|a\|_{L^1}$ is non zero and for $k \in \mathbb{N}$, let $P_k := a|D|^{-k}$. Then $P_k \in OP^{-k} \subset OP^0$ and its associated zeta-function has a pole at $d - k$:

$$\begin{aligned} \operatorname{Res}_{s=d-k} \zeta_{\mathcal{D}}^P(s) &= \operatorname{Res}_{s=0} \zeta_{\mathcal{D}}^P(s + d - k) = \operatorname{Res}_{s=0} \operatorname{Tr} \left(a|D|^{-k} |D|^{-(s+d-k)} \right) = \int a|D|^{-d} \\ &= \int_M a(x) \int_{S_x^* M} \operatorname{Tr} \left((\sigma_1^{|\mathcal{D}|})^{-d}(x, \xi) \right) |d\xi| |dx| = c \int_M a(x) \int_{S_x^* M} \|\xi\|^{-d} |d\xi| |dx| \\ &= c \int_M a(x) d\operatorname{vol}_g(x) = c \|a\|_{L^1} \neq 0. \end{aligned}$$

Conversely, since $\Psi^0(\mathcal{A})$ is contained in the algebra of all pseudodifferential operators of order less or equal to 0, it is known [52, 111, 112] that $Sp(M) \subset \{d - k : k \in \mathbb{N}\}$ as seen in Theorem 4.1.

All poles are simple since \mathcal{D} being differential and M being without boundary, $a'_k = 0$, for all $k \in \mathbb{N}^*$ in (65). \square

Remark 5.35. Due to our efforts to mimic the commutative case, we get as in Theorem 1.22 that the noncommutative integral is a trace on $\Psi^*(\mathcal{A})$. However, when the dimension spectrum is not simple, the analog of $W\operatorname{Res}$ is no longer a trace.

The equation (58) can be obtained via (12) and (66) since $\sigma_d^{|\mathcal{D}_A|^{-d}} = \sigma_d^{|\mathcal{D}|^{-d}}$. In dimension $d = 4$, the computation in (38) of coefficient $a_4(1, \mathcal{D}_A^2)$ gives

$$\zeta_{\mathcal{D}_A}(0) = c_1 \int_M (5R^2 - 8R_{\mu\nu\rho\sigma}R^{\mu\nu\rho\sigma} - 7R_{\mu\nu\rho\sigma}R^{\mu\nu\rho\sigma}) d\text{vol} + c_2 \int_M \text{tr}(F_{\mu\nu}F^{\mu\nu}) d\text{vol},$$

see Corollary 6.4 to see precise correspondence between $a_k(1, \mathcal{D}_A^2)$ and $\zeta_{\mathcal{D}_A}(0)$. One recognizes the Yang–Mills action which will be generalized in Section 6.1.3 to arbitrary spectral triples. According to Corollary 5.31, a commutative geometry has no tadpoles.

5.7 Scalar curvature

What could be the scalar curvature of a spectral triple $(\mathcal{A}, \mathcal{H}, \mathcal{D})$? Of course, we need to consider first the case of a commutative geometry $(C^\infty(M), L^2(M, S), \mathcal{D})$ of dimension $d = 4$: We know that $\int f(x) \mathcal{D}^{-d+2} = \int_M f(x) s(x) d\text{vol}(x)$ where s is the scalar curvature for any $f \in C^\infty(M)$. This suggests the following

Definition 5.36. Let $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ be a spectral triple of dimension d . The scalar curvature is the map $\mathcal{R} : a \in \mathcal{A} \rightarrow \mathbb{C}$ defined by

$$\mathcal{R}(a) := \int a \mathcal{D}^{-d+2}.$$

In the commutative case, \mathcal{R} is a trace on the algebra. More generally

Proposition 5.37. If \mathcal{R} is a trace on \mathcal{A} and the tadpoles $\int A \mathcal{D}^{-d+1}$ are zero for all $A \in \Omega_{\mathcal{D}}^1$, \mathcal{R} is invariant by inner fluctuations $\mathcal{D} \rightarrow \mathcal{D} + A$.

See [31, Proposition 1.153] for a proof.

5.8 Tensor product of spectral triples

There is a natural notion of tensor for spectral triples which corresponds to direct product of manifolds in the commutative case. Let $(\mathcal{A}_i, \mathcal{D}_i, \mathcal{H}_i)$, $i = 1, 2$, two spectral triples of dimension d_i with simple dimension spectrum. Assume the first to be of even dimension, with grading χ_1 .

The spectral triple $(\mathcal{A}, \mathcal{D}, \mathcal{H})$ associated to the tensor product is defined by

$$\mathcal{A} := \mathcal{A}_1 \otimes \mathcal{A}_2, \quad \mathcal{D} := \mathcal{D}_1 \otimes 1 + \chi_1 \otimes \mathcal{D}_2, \quad \mathcal{H} := \mathcal{H}_1 \otimes \mathcal{H}_2.$$

The interest of χ_1 is to guarantee additivity: $\mathcal{D}^2 = \mathcal{D}_1^2 \otimes 1 + 1 \otimes \mathcal{D}_2^2$.

We assume that

$$\text{Tr}(e^{-t\mathcal{D}_1^2}) \sim_{t \rightarrow 0} a_1 t^{-d_1/2}, \quad \text{Tr}(e^{-t\mathcal{D}_2^2}) \sim_{t \rightarrow 0} a_2 t^{-d_2/2}. \quad (68)$$

Lemma 5.38. The triple $(\mathcal{A}, \mathcal{D}, \mathcal{H})$ has dimension $d = d_1 + d_2$.

Moreover, the function $\zeta_{\mathcal{D}}(s) = \text{Tr}(|D|^{-s})$ has a simple pole at $s = d_1 + d_2$ with

$$\text{Res}_{s=d_1+d_2} (\zeta_{\mathcal{D}}(s)) = \frac{1}{2} \frac{\Gamma(d_1/2)\Gamma(d_2/2)}{\Gamma(d/2)} \text{Res}_{s=d_1} (\zeta_{\mathcal{D}_1}(s)) \text{Res}_{s=d_2} (\zeta_{\mathcal{D}_2}(s)).$$

Proof. If $(\mu_n(A))$ are the singular values of A ,

$$\zeta_{\mathcal{D}}(2s) = \sum_{n=0}^{\infty} \mu_n(D_1^2 \otimes 1 + 1 \otimes D_2^2)^{-s} = \sum_{n,m=0}^{\infty} \left(\mu_n(D_1^2) + \mu_m(D_2^2) \right)^{-s}.$$

Since $\left(\mu_n(D_1^2) + \mu_m(D_2^2) \right)^{-(c_1+c_2)/2} \leq \mu_n(D_1)^{-c_1} \mu_m(D_2)^{-c_2}$, this shows in particular that $\zeta_{\mathcal{D}}(c_1 + c_2) \leq \zeta_{\mathcal{D}_1}(c_1) \zeta_{\mathcal{D}_2}(c_2)$ if $c_i > d_i$, and in particular that

$$d := \inf \{ c \in \mathbb{R}^+ : \zeta_{\mathcal{D}}(c) < \infty \} \leq d_1 + d_2.$$

We claim that $d = d_1 + d_2$: recall first that in (68)

$$\begin{aligned} a_i &:= \text{Res}_{s=d_i/2} \left(\Gamma(s) \zeta_{\mathcal{D}_i}(2s) \right) = \Gamma(d_i/2) \text{Res}_{s=d_i/2} \left(\zeta_{\mathcal{D}_i}(2s) \right) \\ &= \frac{1}{2} \Gamma(d_i/2) \text{Res}_{s=d_i} \left(\zeta_{\mathcal{D}_i}(s) \right). \end{aligned} \quad (69)$$

If $f(s) := \Gamma(s) \zeta_{\mathcal{D}}(2s)$,

$$\begin{aligned} f(s) &= \Gamma(s) \text{Tr} \left(\mathcal{D}^{-2s} \right) = \text{Tr} \left(\int_0^{\infty} e^{-t\mathcal{D}^2} t^{s-1} dt \right) = \int_0^1 \text{Tr} \left(e^{-t\mathcal{D}^2} \right) t^{s-1} dt + g(s) \\ &= \int_0^1 \text{Tr} \left(e^{-t\mathcal{D}_1^2} \right) \text{Tr} \left(e^{-t\mathcal{D}_2^2} \right) t^{s-1} dt + g(s) \end{aligned}$$

where g is a holomorphic function since the map $x \in \mathbb{R} \rightarrow \int_1^{\infty} e^{-tx^2} t^{x-1} dt$ is in Schwartz space.

Since $\text{Tr} \left(e^{-t\mathcal{D}_1^2} \right) \text{Tr} \left(e^{-t\mathcal{D}_2^2} \right) \sim_{t \rightarrow 0} a_1 a_2 t^{-(d_1+d_2)/2}$, we get that the function $f(s)$ has a simple pole at $s = (d_1 + d_2)/2$. We conclude that $\zeta_{\mathcal{D}}(s)$ has a simple pole at $s = d_1 + d_2$.

Moreover, thanks to (69),

$$\frac{1}{2} \Gamma((d_1 + d_2)/2) \text{Res}_{s=d} \left(\zeta_{\mathcal{D}}(s) \right) = \frac{1}{2} \Gamma(d_1/2) \text{Res}_{s=d_1} \left(\zeta_{\mathcal{D}_1}(s) \right) \frac{1}{2} \Gamma(d_2/2) \text{Res}_{s=d_2} \left(\zeta_{\mathcal{D}_2}(s) \right). \quad \square$$

6 Spectral action

6.1 On the search for a good action functional

We would like to obtain a good action for any spectral triple and for this it is useful to look at some examples in physics.

In any physical theory based on geometry, the interest of an action functional is, by a minimization process, to exhibit a particular geometry, for instance, trying to distinguish between different metrics. This is the case in general relativity with the Einstein–Hilbert action (with its Riemannian signature).

6.1.1 Einstein–Hilbert action

This action is

$$S_{EH}(g) := - \int_M s_g(x) dvol_g(x) \quad (70)$$

where s is the scalar curvature (chosen positive for the sphere). This is nothing else (up to a constant, in dimension 4) than $f \mathcal{D}^{-2}$ as quoted after (40).

This action is interesting for the following reason: Let \mathcal{M}_1 be the set of Riemannian metrics g on M such that $\int_M dvol_g = 1$. By a theorem of Hilbert [5], $g \in \mathcal{M}_1$ is a critical point of $S_{EH}(g)$ restricted to \mathcal{M}_1 if and only if (M, g) is an Einstein manifold (the Ricci curvature R of g is proportional by a constant to g : $R = c g$). Taking the trace, this means that $s_g = c \dim(M)$ and such manifold have a constant scalar curvature.

But in the search for invariants under diffeomorphisms, they are more quantities than the Einstein–Hilbert action, a trivial example being $\int_M f(s_g(x)) dvol_g(x)$ and they are others [43]. In this desire to implement gravity in noncommutative geometry, the eigenvalues of the Dirac operator look as natural variables [70]. However we are looking for observables which add up under disjoint unions of different geometries.

6.1.2 Quantum approach and spectral action

In a way, a spectral triple fits quantum field theory since \mathcal{D}^{-1} can be seen as the propagator (or line element ds) for (Euclidean) fermions and we can compute Feynman graphs with fermionic internal lines. As explained in section 5.4, the gauge bosons are only derived objects obtained from internal fluctuations via Morita equivalence given by a choice of a connection which is associated to a one-form in $\Omega_{\mathcal{D}}^1(\mathcal{A})$. Thus, the guiding principle followed by Connes and Chamseddine is to use a theory which is pure gravity with a functional action based on the spectral triple, namely which depends on the spectrum of \mathcal{D} [11]. They proposed the following

Definition 6.1. *The spectral action of a spectral triple $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ is defined by*

$$\mathcal{S}(\mathcal{D}, f, \Lambda) := \text{Tr} \left(f(\mathcal{D}^2/\Lambda^2) \right)$$

where $\Lambda \in \mathbb{R}^+$ plays the role of a cut-off and f is any positive function (such that $f(\mathcal{D}^2/\Lambda^2)$ is a trace-class operator).

Remark 6.2. We can also define $\mathcal{S}(\mathcal{D}, f, \Lambda) = \text{Tr} \left(f(\mathcal{D}/\Lambda) \right)$ when f is positive and even. With this second definition, $S(\mathcal{D}, g, \Lambda) = \text{Tr} \left(f(\mathcal{D}^2/\Lambda^2) \right)$ with $g(x) := f(x^2)$.

For f , one can think of the characteristic function of $[-1, 1]$, thus $f(\mathcal{D}/\Lambda)$ is nothing else but the number of eigenvalues of \mathcal{D} within $[-\Lambda, \Lambda]$.

When this action has an asymptotic series in $\Lambda \rightarrow \infty$, we deal with an effective theory. Naturally, \mathcal{D} has to be replaced by \mathcal{D}_A which is just a decoration. To this bosonic part of the action, one adds a fermionic term $\frac{1}{2} \langle J\psi, \mathcal{D}\psi \rangle$ for $\psi \in \mathcal{H}$ to get a full action. In the standard model of particle physics, this latter corresponds to the integration of the Lagrangian part for the coupling between gauge bosons and Higgs bosons with fermions. Actually, the finite dimension part of the noncommutative standard model is of KO -dimension 6, thus $\langle \psi, \mathcal{D}\psi \rangle$ has to be replaced by $\frac{1}{2} \langle J\psi, \mathcal{D}\psi \rangle$ for $\psi = \chi\psi \in \mathcal{H}$, see [31].

6.1.3 Yang–Mills action

This action plays an important role in physics so it is natural to consider it in the noncommutative framework. Recall first the classical situation: let G be a compact Lie group with its Lie algebra \mathfrak{g} and let $A \in \Omega^1(M, \mathfrak{g})$ be a connection. If $F := da + \frac{1}{2}[A, A] \in \Omega^2(M, \mathfrak{g})$ is the curvature (or field strength) of A , then the Yang-Mills action is $S_{YM}(A) = \int_M \text{tr}(F \wedge \star F) \text{dvol}_g$. In the abelian case $G = U(1)$, it is the Maxwell action and its quantum version is the quantum electrodynamics (QED) since the un-gauged $U(1)$ of electric charge conservation can be gauged and its gauging produces electromagnetism [99]. It is conformally invariant when the dimension of M is $d = 4$.

The study of its minima and its critical values can also be made for a spectral triple $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ of dimension d [24, 25]: let $A \in \Omega_{\mathcal{D}}^1(\mathcal{A})$ and curvature $\theta = dA + A^2$; then it is natural to consider

$$I(A) := \text{Tr}_{Dix}(\theta^2 |\mathcal{D}|^{-d})$$

since it coincides (up to a constant) with the previous Yang-Mills action in the commutative case: if $P = \theta^2 |\mathcal{D}|^{-d}$, then Theorems 1.22 and 2.14 give the claim since for the principal symbol, $\text{tr}(\sigma^P(x, \xi)) = c \text{tr}(F \wedge \star F)(x)$.

There is nevertheless a problem with the definition of dA : if $A = \sum_j \pi(a_j)[\mathcal{D}, \pi(b_j)]$, then $dA = \sum_j [\mathcal{D}, \pi(a_j)][\mathcal{D}, \pi(b_j)]$ can be non-zero while $A = 0$. This ambiguity means that, to get a graded differential algebra $\Omega_{\mathcal{D}}^*(\mathcal{A})$, one must divide by a junk, for instance $\Omega_{\mathcal{D}}^2 \simeq \pi(\Omega^2/\pi(\delta(\text{Ker}(\pi) \cap \Omega^1)))$ where $\Omega^k(\mathcal{A})$ is the set of universal k -forms over \mathcal{A} given by the set of $a_0 \delta a_1 \cdots \delta a_k$ (before representation on \mathcal{H} : $\pi(a_0 \delta a_1 \cdots \delta a_k) := a_0 [\mathcal{D}, a_1] \cdots [\mathcal{D}, a_k]$). Let \mathcal{H}_k be the Hilbert space completion of $\pi(\Omega^k(\mathcal{A}))$ with the scalar product defined by $\langle A_1, A_2 \rangle_k := \text{Tr}_{Dix}(A_2^* A_1 |\mathcal{D}|^{-d})$ for $A_j \in \pi(\Omega^k(\mathcal{A}))$.

The Yang–Mills action on $\Omega^1(\mathcal{A})$ is

$$S_{YM}(V) := \langle \delta V + V^2, \delta V + V^2 \rangle. \quad (71)$$

It is positive, quartic and gauge invariant under $V \rightarrow \pi(u)V\pi(u^*) + \pi(u)[\mathcal{D}, \pi(u^*)]$ when $u \in \mathcal{U}(\mathcal{A})$. Moreover,

$$S_{YM}(V) = \inf \{ I(\omega) \mid \omega \in \Omega^1(\mathcal{A}), \pi(\omega) = V \}$$

since the above ambiguity disappears when taking the infimum.

This Yang–Mills action can be extended to the equivalent of Hermitean vector bundles on M , namely finitely projective modules over \mathcal{A} .

The spectral action is more conceptual than the Yang–Mills action since it gives no fundamental role to the distinction between gravity and matter in the artificial decomposition $\mathcal{D}_A = \mathcal{D} + A$. For instance, for the minimally coupled standard model, the Yang–Mills action for the vector potential is part of the spectral action, as far as the Einstein–Hilbert action for the Riemannian metric [12].

As quoted in [17], the spectral action has conceptual advantages:

- Simplicity: when f is a cutoff function, the spectral action is just the counting function.
- Positivity: when f is positive (which is the case for a cutoff function), the action $\text{Tr} \left(f(\mathcal{D}/\Lambda) \right) \geq 0$ has the correct sign for a Euclidean action: the positivity of the function f will insure that the actions for gravity, Yang–Mills, Higgs couplings are all positive and the Higgs mass term is negative.

- Invariance: the spectral action has a much stronger invariance group than the usual diffeomorphism group as for the gravitational action; this is the unitary group of the Hilbert space \mathcal{H} .

However, this action is not local. It only becomes so when it is replaced by the asymptotic expansion:

6.2 Asymptotic expansion for $\Lambda \rightarrow \infty$

The heat kernel method already used in previous sections will give a control of spectral action $S(\mathcal{D}, f, \Lambda)$ when Λ goes to infinity.

Theorem 6.3. *Let $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ be a spectral triple with a simple dimension spectrum Sd . We assume that*

$$\text{Tr} \left(e^{-t\mathcal{D}^2} \right) \underset{t \downarrow 0}{\sim} \sum_{\alpha \in Sd} a_\alpha t^\alpha \quad \text{with } a_\alpha \neq 0. \quad (72)$$

Then, for the zeta function $\zeta_{\mathcal{D}}$ defined in (44)

$$a_\alpha = \frac{1}{2} \text{Res}_{s=-2\alpha} \left(\Gamma(s/2) \zeta_{\mathcal{D}}(s) \right). \quad (73)$$

(i) If $\alpha < 0$, $\zeta_{\mathcal{D}}$ has a pole at -2α with $a_\alpha = \frac{1}{2} \Gamma(-\alpha) \text{Res}_{s=-2\alpha} \zeta_{\mathcal{D}}(s)$.

(ii) For $\alpha = 0$, we get $a_0 = \zeta_{\mathcal{D}}(0) + \dim \text{Ker } \mathcal{D}$.

(iii) If $\alpha > 0$, $a_\alpha = \zeta(-2\alpha) \text{Res}_{s=-\alpha} \Gamma(s)$.

(iv) The spectral action has the asymptotic expansion over the positive part Sd^+ of Sd :

$$\text{Tr} \left(f(\mathcal{D}/\Lambda) \right) \underset{\Lambda \rightarrow +\infty}{\sim} \sum_{\beta \in Sd^+} f_\beta \Lambda^\beta \int |\mathcal{D}|^\beta + f(0) \zeta_{\mathcal{D}}(0) + \dots \quad (74)$$

where the dependence of the even function f is $f_\beta := \int_0^\infty f(x) x^{\beta-1} dx$ and \dots involves the full Taylor expansion of f at 0.

Proof. (i): Since $\Gamma(s/2) |\mathcal{D}|^{-s} = \int_0^\infty e^{-t\mathcal{D}^2} t^{s/2-1} dt = \int_0^1 e^{-t\mathcal{D}^2} t^{s/2-1} dt + f(s)$, where the function f is holomorphic (since the map $x \rightarrow \int_1^\infty e^{-tx^2} x^{s/2-1} dt$ is in the Schwartz space), the swap of $\text{Tr} \left(e^{-t\mathcal{D}^2} \right)$ with a sum of $a_\alpha t^\alpha$ and $a_\alpha \int_0^1 t^{\alpha+s/2-1} dt = \frac{2a_\alpha}{s+2\alpha}$ yields (73).

(ii): The regularity of $\Gamma(s/2)^{-1} \underset{s \rightarrow 0}{\sim} s/2$ around zero implies that only the pole part at $s = 0$ of $\int_0^\infty \text{Tr} \left(e^{-t\mathcal{D}^2} \right) t^{s/2-1} dt$ contributes to $\zeta_{\mathcal{D}}(0)$. This contribution is $a_0 \int_0^1 t^{s/2-1} dt = \frac{2a_0}{s}$.

(iii) follows from (73).

(iv): Assume $f(x) = g(x^2)$ where g is a Laplace transform: $g(x) := \int_0^\infty e^{-sx} \phi(s) ds$. We will see in Section 6.3 how to relax this hypothesis. Since $g(t\mathcal{D}^2) = \int_0^\infty e^{-st\mathcal{D}^2} \phi(s) ds$, $\text{Tr} \left(g(t\mathcal{D}^2) \right) \underset{t \downarrow 0}{\sim} \sum_{\alpha \in \text{Sp}^+} a_\alpha t^\alpha \int_0^\infty s^\alpha g(s) ds$. When $\alpha < 0$, $s^\alpha = \Gamma(-\alpha)^{-1} \int_0^\infty e^{-sy} y^{-\alpha-1} dy$ and $\int_0^\infty s^\alpha \phi(s) ds = \Gamma(-\alpha)^{-1} \int_0^\infty g(y) y^{-\alpha-1} dy$. Thus

$$\text{Tr} \left(g(t\mathcal{D}^2) \right) \underset{t \downarrow 0}{\sim} \sum_{\alpha \in \text{Sp}^-} \left[\frac{1}{2} \text{Res}_{s=-2\alpha} \zeta_{\mathcal{D}}(s) \int_0^\infty g(y) y^{-\alpha-1} dy \right] t^\alpha.$$

Thus (74) follows from (i), (ii) and $\frac{1}{2} \int_0^\infty g(y) y^{\beta/2-1} dy = \int_0^\infty f(x) x^{\beta-1} dx$. \square

It can be useful to make a connection with (39) of Section 4.3:

Corollary 6.4. *Assume that the spectral triple $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ has dimension d . If*

$$\text{Tr} \left(e^{-t\mathcal{D}^2} \right) \underset{t \downarrow 0}{\sim} \sum_{k \in \{0, \dots, d\}} t^{(k-d)/2} a_k(\mathcal{D}^2) + \dots, \quad (75)$$

then

$$\mathcal{S}(\mathcal{D}, f, \Lambda) \underset{t \downarrow 0}{\sim} \sum_{k \in \{1, \dots, d\}} f_k \Lambda^k a_{d-k}(\mathcal{D}^2) + f(0) a_d(\mathcal{D}^2) + \dots$$

with $f_k := \frac{1}{\Gamma(k/2)} \int_0^\infty f(s) s^{k/2-1} ds$.

Moreover,

$$\begin{aligned} a_k(\mathcal{D}^2) &= \frac{1}{2} \Gamma\left(\frac{d-k}{2}\right) \int |\mathcal{D}|^{-d+k} \text{ for } k = 0, \dots, d-1, \\ a_d(\mathcal{D}^2) &= \dim \text{Ker } \mathcal{D} + \zeta_{\mathcal{D}^2}(0). \end{aligned} \quad (76)$$

Proof. We rewrite the hypothesis on $\text{Tr} \left(e^{-t\mathcal{D}^2} \right)$ as

$$\text{Tr} \left(e^{-t\mathcal{D}^2} \right) \underset{t \downarrow 0}{\sim} \sum_{\alpha \in \{-d/2, \dots, -1/2\}} A_\alpha t^\alpha + A_0 = \sum_{k \in \{1, \dots, d\}} t^{(k-d)/2} a_k(\mathcal{D}^2) + a_d(\mathcal{D}^2)$$

with $a_k(\mathcal{D}^2) := A_{(k-d)/2}$.

For $\alpha < 0$, we repeat the above proof:

$$\begin{aligned} \mathcal{S}(\mathcal{D}, f, \Lambda) &\underset{t \downarrow 0}{\sim} \sum_{\alpha \in \{-d/2, \dots, -1/2\}} A_\alpha \Lambda^{-2\alpha} \frac{1}{\Gamma(-\alpha)} \int_0^\infty f(s) s^{-\alpha-1} ds + A_0 \\ &= \sum_{l \in \{1, \dots, d\}} A_{(l-d)/2} \Lambda^{d-l} \frac{1}{\Gamma((d-l)/2)} \int_0^\infty f(s) s^{(d-l)/2-1} ds + A_0 f(0) \\ &= \sum_{l \in \{1, \dots, d\}} a_l \Lambda^{d-l} \frac{1}{\Gamma((d-l)/2)} \int_0^\infty f(s) s^{(d-l)/2-1} ds + a_d f(0) \\ &= \sum_{k \in \{1, \dots, d\}} a_{d-k} \Lambda^k \frac{1}{\Gamma(k/2)} \int_0^\infty f(s) s^{k/2-1} ds + a_d f(0). \end{aligned}$$

Again, for $\alpha < 0$,

$$\begin{aligned} A_\alpha &= \frac{1}{2}\Gamma(-\alpha) \operatorname{Res}_{s=-2\alpha} \operatorname{Tr} \left(|\mathcal{D}|^{-s} \right) = \frac{1}{2}\Gamma(-\alpha) \operatorname{Res}_{s=0} \operatorname{Tr} \left(|\mathcal{D}|^{-(s-2\alpha)} \right) \\ &= \frac{1}{2}\Gamma(-\alpha) \operatorname{Res}_{s=0} \operatorname{Tr} \left(|\mathcal{D}|^{2\alpha} |\mathcal{D}|^{-s} \right) = \frac{1}{2}\Gamma(-\alpha) \int |\mathcal{D}|^{2\alpha}. \end{aligned}$$

Thus, for $\alpha = \frac{k-d}{2} < 0$, $k = 0, \dots, d-1$,

$$a_k(\mathcal{D}^2) = A_{(k-d)/2} = \frac{1}{2}\Gamma\left(\frac{d-k}{2}\right) \int |\mathcal{D}|^{-d+k}. \quad \square$$

The asymptotics (74) uses the value of $\zeta_{\mathcal{D}}(0)$ in the constant term Λ^0 , so it is fundamental to look at its variation under a gauge fluctuation $\mathcal{D} \rightarrow \mathcal{D} + A$ as we saw in Theorem 5.17.

6.3 Remark on the use of Laplace transform

The spectral action asymptotic behavior

$$S(\mathcal{D}, f, \Lambda) \underset{\Lambda \rightarrow +\infty}{\sim} \sum_{n=0}^{\infty} c_n \Lambda^{d-n} a_n(\mathcal{D}^2) \quad (77)$$

has been proved for a smooth function f which is a Laplace transform for an arbitrary spectral triple (with simple dimension spectrum) satisfying (72). However, this hypothesis is too restrictive since it does not cover the heat kernel case where $f(x) = e^{-x}$.

When the triple is commutative and \mathcal{D}^2 is a generalized Laplacian on sections of a vector bundle over a manifold of dimension 4, Estrada–Gracia-Bondía–Várilly proved in [38] that previous asymptotics is

$$\begin{aligned} \operatorname{Tr} \left(f(\mathcal{D}^2/\Lambda^2) \right) &\sim \frac{1}{(4\pi)^2} \left[\left(\operatorname{rk}(\mathbb{E}) \int_0^\infty x f(x) dx \right) \Lambda^4 + \left(b_2(\mathcal{D}^2) \int_0^\infty f(x) dx \right) \Lambda^2 \right. \\ &\quad \left. + \left(\sum_{m=0}^{\infty} (-1)^m f^{(m)}(0) b_{2m+4}(\mathcal{D}^2) \right) \Lambda^{-2m} \right], \quad \Lambda \rightarrow \infty \end{aligned}$$

where $(-1)^m b_{2m+4}(\mathcal{D}^2) = \frac{(4\pi)^2}{m!} \mu_m(\mathcal{D}^2)$ are suitably normalized, integrated moment terms of the spectral density of \mathcal{D}^2 .

The main point is that *this asymptotics makes sense in the Cesàro sense* (see [38] for definition) for $f \in \mathcal{K}'(\mathbb{R})$, which is the dual of $\mathcal{K}(\mathbb{R})$. This latter is the space of smooth functions ϕ such that for some $a \in \mathbb{R}$, $\phi^{(k)}(x) = \mathcal{O}(|x|^{a-k})$ as $|x| \rightarrow \infty$, for each $k \in \mathbb{N}$. In particular, the Schwartz functions are in $\mathcal{K}(\mathbb{R})$ (and even dense).

Of course, the counting function is not smooth but is in $\mathcal{K}'(\mathbb{R})$, so such behavior (77) is wrong beyond the first term, but is correct in the Cesàro sense. Actually there are more derivatives of f at 0 as explained on examples in [38, p. 243]. See also Section 9.4.

6.4 About convergence and divergence, local and global aspects of the asymptotic expansion

The asymptotic expansion series (75) of the spectral action may or may not converge. It is known that each function $g(\Lambda^{-1})$ defines at most a unique expansion series when $\Lambda \rightarrow \infty$

but the converse is not true since several functions have the same asymptotic series. We give here examples of convergent and divergent series of this kind.

When M is the torus \mathbb{T}^d as in Example 2.13 with $\Delta = \delta^{\mu\nu} \partial_\mu \partial_\nu$,

$$\mathrm{Tr}(e^{t\Delta}) = \frac{(4\pi)^{-d/2} \mathrm{Vol}(\mathbb{T}^d)}{t^{d/2}} + \mathcal{O}(t^{-d/2} e^{-1/4t}),$$

thus the asymptotic series $\mathrm{Tr}(e^{t\Delta}) \simeq \frac{(4\pi)^{-d/2} \mathrm{Vol}(\mathbb{T}^d)}{t^{d/2}}$, $t \rightarrow 0$, has only one term.

In the opposite direction, let now M be the unit four-sphere \mathbb{S}^4 and \mathcal{D} be the usual Dirac operator. By Proposition 5.34, equation (72) yields (see [20]):

$$\begin{aligned} \mathrm{Tr}(e^{-t\mathcal{D}^2}) &= \frac{1}{t^2} \left(\frac{2}{3} + \frac{2}{3}t + \sum_{k=0}^n a_k t^{k+2} + \mathcal{O}(t^{n+3}) \right), \\ a_k &:= \frac{(-1)^k 4}{3 k!} \left(\frac{B_{2k+2}}{2k+2} - \frac{B_{2k+4}}{2k+4} \right) \end{aligned}$$

with Bernoulli numbers B_{2k} . Thus $t^2 \mathrm{Tr}(e^{-t\mathcal{D}^2}) \simeq \frac{2}{3} + \frac{2}{3}t + \sum_{k=0}^\infty a_k t^{k+2}$ when $t \rightarrow 0$ and this series is not convergent but only asymptotic:

$$a_k > \frac{4}{3 k!} \frac{|B_{2k+4}|}{2k+4} > 0 \text{ and } |B_{2k+4}| = 2 \frac{(2k+4)!}{(2\pi)^{2k+4}} \zeta(2k+4) \simeq 4\sqrt{\pi(k+2)} \left(\frac{k+2}{\pi e}\right)^{2k+4} \rightarrow \infty \text{ if } k \rightarrow \infty.$$

More generally, in the commutative case considered above and when \mathcal{D} is a differential operator—like a Dirac operator, the coefficients of the asymptotic series of $\mathrm{Tr}(e^{-t\mathcal{D}^2})$ are locally defined by the symbol of \mathcal{D}^2 at point $x \in M$ but this is not true in general: in [46] is given a positive elliptic pseudodifferential such that non-locally computable coefficients especially appear in (75) when $2k > d$. Nevertheless, all coefficients are local for $2k \leq d$.

Recall that a locally computable quantity is the integral on the manifold of a local frame-independent smooth function of one variable, depending only on a finite number of derivatives of a finite number of terms in the asymptotic expansion of the total symbol of \mathcal{D}^2 . For instance, some nonlocal information contained in the ultraviolet asymptotics can be recovered if one looks at the (integral) kernel of $e^{-t\sqrt{-\Delta}}$: in \mathbb{T}^1 , with $\mathrm{Vol}(\mathbb{T}^1) = 2\pi$, we get [39]

$$\mathrm{Tr}(e^{-t\sqrt{-\Delta}}) = \frac{\sinh(t)}{\cosh(t) - 1} = \coth\left(\frac{t}{2}\right) = \frac{2}{t} \sum_{k=0}^\infty \frac{B_{2k}}{(2k)!} t^{2k} = \frac{2}{t} \left[1 + \frac{t^2}{12} - \frac{t^4}{720} + \mathcal{O}(t^6) \right]$$

and the series converges when $t < 2\pi$, since $\frac{B_{2k}}{(2k)!} = (-1)^{k+1} \frac{2\zeta(2k)}{(2\pi)^{2k}}$, thus $\frac{|B_{2k}|}{(2k)!} \simeq \frac{2}{(2\pi)^{2k}}$ when $k \rightarrow \infty$.

Thus we have an example where $t \rightarrow \infty$ cannot be used with the asymptotic series.

Thus the spectral action of Corollary 6.4 precisely encodes these local and nonlocal behavior which appear or not in its asymptotics for different f . The coefficient of the action for the positive part (at least) of the dimension spectrum correspond to renormalized traces, namely the noncommutative integrals of (76). In conclusion, the asymptotic (77) of spectral action may or may not have nonlocal coefficients.

For the flat torus \mathbb{T}^d , the difference between $\mathrm{Tr}(e^{t\Delta})$ and its asymptotic series is a term which is related to periodic orbits of the geodesic flow on \mathbb{T}^d . Similarly, the counting function $N(\lambda)$ (number of eigenvalues including multiplicities of Δ less than λ) obeys Weyl's law: $N(\lambda) = \frac{(4\pi)^{-d/2} \mathrm{Vol}(\mathbb{T}^d)}{\Gamma(d/2+1)} \lambda^{d/2} + o(\lambda^{d/2})$ — see [1] for a nice historical review on these fundamental points. The relationship between the asymptotic expansion of the heat kernel and the formal expansion of the spectral measure is clear: the small- t asymptotics of heat kernel

is determined by the large- λ asymptotics of the density of eigenvalues (and eigenvectors). However, the latter is defined modulo some average: Cesàro sense as reminded in Section 6.3, or Riesz mean of the measure which washes out ultraviolet oscillations, but also gives informations on intermediate values of λ [39].

In [17, 77] are given examples of spectral actions on (compact) commutative geometries of dimension 4 whose asymptotics have only two terms. In the quantum group $SU_q(2)$, the spectral action itself has only 4 terms, independently of the choice of function f .

See [63] for more examples.

6.5 About the physical meaning of the spectral action via its asymptotics

As explain before, the spectral action is non-local. Its localization does not cover all situations: consider for instance the commutative geometry of a spin manifold M (see Section 5.6) of dimension 4. One adds a gauge connection $A \in \Gamma^\infty(M, \text{End}(S))$ to the Dirac operator \mathcal{D} such that $\mathcal{D} = i\gamma^\mu(\partial_\mu + A_\mu)$, thus with a field strength $F_{\mu\nu} = \partial_\mu A_\nu - \partial_\nu A_\mu + [A_\mu, A_\nu]$. We can apply (32) with $P = \mathcal{D}^2$ and find the 3 coefficients $a_i(1, P)$ of (38) with $i = 0, 2, 4$. The expansion (77) corresponds to a weak field expansion.

Moreover a commutative geometry times a finite one where the finite one is algebra is a sum of matrices (like in Remark 5.33) has been deeply and intensively investigated for the noncommutative approach to standard model of particle physics, see [21, 31]. This approach offers a lot of interesting perspectives, for instance, the possibility to compute the Higgs representations and mass (for each noncommutative model) is particularly instructive [11, 16, 18, 58, 64, 65, 72, 80]. Of course, since the first term in (74) is a cosmological term, one may be worried by its large value (for instance in the noncommutative standard model where the cutoff is, roughly speaking the Planck scale). At the classical level, one can work with unimodular gravity where the metric (so the Dirac operator) \mathcal{D} varies within the set \mathcal{M}_1 of metrics which preserve the volume as in Section 6.1.1. Thus it remains only (!) to control the inflaton: see [14].

The spectral action has been computed in [61] for the quantum group $SU_q(2)$ which is not a deformation of $SU(2)$ of the type considered in Section 9.5 on the Moyal plane. It is quite peculiar since (74) has only a finite number of terms.

Due to the difficulties to deal with non-compact manifolds (see nevertheless Section 9), the case of spheres \mathbb{S}^4 or $\mathbb{S}^3 \times \mathbb{S}^1$ has been investigated in [17, 20] for instance in the case of Robertson–Walker metrics.

All the machinery of spectral geometry as been recently applied to cosmology, computing the spectral action in few cosmological models related to inflation, see [67, 77–79, 83, 97].

Spectral triples associated to manifolds with boundary have been considered in [15, 19, 19, 59, 60, 62]. The main difficulty is precisely to put nice boundary conditions to the operator \mathcal{D} to still get a selfadjoint operator and then, to define a compatible algebra \mathcal{A} . This is probably a must to obtain a result in a noncommutative Hamiltonian theory in dimension 1+3.

The case of manifolds with torsion has also been studied in [54, 86, 87], and even with boundary in [62]. These works show that the Holst action appears in spectral actions and that torsion could be detected in a noncommutative world.

7 Residues of series and integral, holomorphic continuation, etc

The aim of this section is to control the holomorphy of series of holomorphic functions. The necessity of a Diophantine condition appears quite naturally. This section has its own interest, but will be fully applied in the next one devoted to the noncommutative torus. The main idea is to get a condition which guarantee the commutation of a residue and a series.

This section is quite technical, but with only non-difficult notions. Nevertheless, the devil is hidden into the details and I recommend to the reader to have a look at the proofs despite their lengths.

Reference: [37].

Notations:

In the following, the prime in \sum' means that we omit terms with division by zero in the summand. B^n (resp. S^{n-1}) is the closed ball (resp. the sphere) of \mathbb{R}^n with center 0 and radius 1 and the Lebesgue measure on S^{n-1} will be noted dS .

For any $x = (x_1, \dots, x_n) \in \mathbb{R}^n$ we denote by $|x| = \sqrt{x_1^2 + \dots + x_n^2}$ the Euclidean norm and $|x|_1 := |x_1| + \dots + |x_n|$.

By $f(x, y) \ll_y g(x)$ uniformly in x , we mean that $|f(x, y)| \leq a(y) |g(x)|$ for all x and y for some $a(y) > 0$.

7.1 Residues of series and integral

In order to be able to compute later the residues of certain series, we prove here the following

Theorem 7.1. *Let $P(X) = \sum_{j=0}^d P_j(X) \in \mathbb{C}[X_1, \dots, X_n]$ be a polynomial function where P_j is the homogeneous part of P of degree j . The function*

$$\zeta^P(s) := \sum'_{k \in \mathbb{Z}^n} \frac{P(k)}{|k|^s}, \quad s \in \mathbb{C}$$

has a meromorphic continuation to the whole complex plane \mathbb{C} .

Moreover $\zeta^P(s)$ is not entire if and only if $\mathcal{P}_P := \{j \mid \int_{u \in S^{n-1}} P_j(u) dS(u) \neq 0\} \neq \emptyset$. In that case, ζ^P has only simple poles at the points $j+n$, $j \in \mathcal{P}_P$, with

$$\operatorname{Res}_{s=j+n} \zeta^P(s) = \int_{u \in S^{n-1}} P_j(u) dS(u).$$

The proof of this theorem is based on the following lemmas.

Lemma 7.2. *For any polynomial $P \in \mathbb{C}[X_1, \dots, X_n]$ of total degree $\delta(P) := \sum_{i=1}^n \deg_{X_i} P$ and any $\alpha \in \mathbb{N}_0^n$, we have*

$$\partial^\alpha \left(P(x) |x|^{-s} \right) \ll_{P, \alpha, n} (1 + |s|)^{|\alpha|_1} |x|^{-\sigma - |\alpha|_1 + \delta(P)}$$

uniformly in $x \in \mathbb{R}^n$, $|x| \geq 1$, where $\sigma = \Re(s)$.

Proof. By linearity, we may assume without loss of generality that $P(X) = X^\gamma$ is a monomial. It is easy to prove (for example by induction on $|\alpha|_1$) that for all $\alpha \in \mathbb{N}_0^n$ and $x \in \mathbb{R}^n \setminus \{0\}$:

$$\partial^\alpha \left(|x|^{-s} \right) = \alpha! \sum_{\substack{\beta, \mu \in \mathbb{N}_0^n \\ \beta + 2\mu = \alpha}} \binom{-s/2}{|\beta|_1 + |\mu|_1} \frac{(|\beta|_1 + |\mu|_1)!}{\beta! \mu!} \frac{x^\beta}{|x|^{\sigma + 2(|\beta|_1 + |\mu|_1)}}.$$

It follows that for all $\alpha \in \mathbb{N}_0^n$, we have uniformly in $x \in \mathbb{R}^n$, $|x| \geq 1$:

$$\partial^\alpha \left(|x|^{-s} \right) \ll_{\alpha, n} (1 + |s|)^{|\alpha|_1} |x|^{-\sigma - |\alpha|_1}. \quad (78)$$

By Leibniz formula and (78), we have uniformly in $x \in \mathbb{R}^n$, $|x| \geq 1$:

$$\begin{aligned} \partial^\alpha \left(x^\gamma |x|^{-s} \right) &= \sum_{\beta \leq \alpha} \binom{\alpha}{\beta} \partial^\beta (x^\gamma) \partial^{\alpha - \beta} \left(|x|^{-s} \right) \\ &\ll_{\gamma, \alpha, n} \sum_{\beta \leq \alpha; \beta \leq \gamma} x^{\gamma - \beta} (1 + |s|)^{|\alpha|_1 - |\beta|_1} |x|^{-\sigma - |\alpha|_1 + |\beta|_1} \\ &\ll_{\gamma, \alpha, n} (1 + |s|)^{|\alpha|_1} |x|^{-\sigma - |\alpha|_1 + |\gamma|_1}. \quad \square \end{aligned}$$

Lemma 7.3. *Let $P \in \mathbb{C}[X_1, \dots, X_n]$ be a polynomial of degree d . Then, the difference*

$$\Delta_P(s) := \sum'_{k \in \mathbb{Z}^n} \frac{P(k)}{|k|^s} - \int_{\mathbb{R}^n \setminus B^n} \frac{P(x)}{|x|^s} dx$$

which is defined for $\Re(s) > d + n$, extends holomorphically on the whole complex plane \mathbb{C} .

Proof. We fix in the sequel a function $\psi \in C^\infty(\mathbb{R}^n, \mathbb{R})$ such that for all $x \in \mathbb{R}^n$

$$0 \leq \psi(x) \leq 1, \quad \psi(x) = 1 \text{ if } |x| \geq 1 \quad \text{and} \quad \psi(x) = 0 \text{ if } |x| \leq 1/2.$$

The function $f(x, s) := \psi(x) P(x) |x|^{-s}$, $x \in \mathbb{R}^n$ and $s \in \mathbb{C}$, is in $\mathcal{C}^\infty(\mathbb{R}^n \times \mathbb{C})$ and depends holomorphically on s .

Lemma 7.2 above shows that f is a ‘‘gauged symbol’’ in the terminology of [53, p. 4]. Thus [53, Theorem 2.1] implies that $\Delta_P(s)$ extends holomorphically on the whole complex plane \mathbb{C} . However, to be complete, we will give here a short proof of Lemma 7.3:

It follows from the classical Euler–Maclaurin formula that for any function $h : \mathbb{R} \rightarrow \mathbb{C}$ of class \mathcal{C}^{N+1} satisfying $\lim_{|t| \rightarrow +\infty} h^{(k)}(t) = 0$ and $\int_{\mathbb{R}} |h^{(k)}(t)| dt < +\infty$ for any $k = 0 \dots, N + 1$, that we have

$$\sum_{k \in \mathbb{Z}} h(k) = \int_{\mathbb{R}} h(t) dt + \frac{(-1)^N}{(N+1)!} \int_{\mathbb{R}} B_{N+1}(t) h^{(N+1)}(t) dt$$

where B_{N+1} is the Bernoulli function of order $N + 1$ (it is a bounded periodic function.)

Fix $m' \in \mathbb{Z}^{n-1}$ and $s \in \mathbb{C}$. Applying to the function $h(t) := \psi(m', t) P(m', t) |(m', t)|^{-s}$ (we use Lemma 7.2 to verify hypothesis), we obtain that for any $N \in \mathbb{N}_0$:

$$\sum_{m_n \in \mathbb{Z}} \psi(m', m_n) P(m', m_n) |(m', m_n)|^{-s} = \int_{\mathbb{R}} \psi(m', t) P(m', t) |(m', t)|^{-s} dt + \mathcal{R}_N(m'; s) \quad (79)$$

where $\mathcal{R}_N(m'; s) := \frac{(-1)^N}{(N+1)!} \int_{\mathbb{R}} B_{N+1}(t) \frac{\partial^{N+1}}{\partial x_n^{N+1}} \left(\psi(m', t) P(m', t) |(m', t)|^{-s} \right) dt$.

By Lemma 7.2,

$$\int_{\mathbb{R}} \left| B_{N+1}(t) \frac{\partial^{N+1}}{\partial x_n^{N+1}} \left(\psi(m', t) P(m', t) |(m', t)|^{-s} \right) \right| dt \ll_{P, n, N} (1 + |s|)^{N+1} (|m'| + 1)^{-\sigma - N + \delta(P)}.$$

Thus $\sum_{m' \in \mathbb{Z}^{n-1}} \mathcal{R}_N(m'; s)$ converges absolutely and define a holomorphic function in the half plane $\{\sigma = \Re(s) > \delta(P) + n - N\}$.

Since N is an arbitrary integer, by letting $N \rightarrow \infty$ and using (79) above, we conclude that:

$$s \mapsto \sum_{(m', m_n) \in \mathbb{Z}^{n-1} \times \mathbb{Z}} \psi(m', m_n) P(m', m_n) |(m', m_n)|^{-s} - \sum_{m' \in \mathbb{Z}^{n-1}} \int_{\mathbb{R}} \psi(m', t) P(m', t) |(m', t)|^{-s} dt$$

has a holomorphic continuation to the whole complex plane \mathbb{C} .

After n iterations, we obtain that

$$s \mapsto \sum_{m \in \mathbb{Z}^n} \psi(m) P(m) |m|^{-s} - \int_{\mathbb{R}^n} \psi(x) P(x) |x|^{-s} dx$$

has a holomorphic continuation to the whole \mathbb{C} .

To finish the proof of Lemma 7.3, it is enough to notice that:

- $\psi(0) = 0$ and $\psi(m) = 1, \forall m \in \mathbb{Z}^n \setminus \{0\}$;
- $s \mapsto \int_{B^n} \psi(x) P(x) |x|^{-s} dx = \int_{\{x \in \mathbb{R}^n: 1/2 \leq |x| \leq 1\}} \psi(x) P(x) |x|^{-s} dx$ is a holomorphic function on \mathbb{C} . □

Proof of Theorem 7.1. Using the polar decomposition of the volume form $dx = \rho^{n-1} d\rho dS$ in \mathbb{R}^n , we get for $\Re(s) > d + n$,

$$\int_{\mathbb{R}^n \setminus B^n} \frac{P_j(x)}{|x|^s} dx = \int_1^\infty \frac{\rho^{j+n-1}}{\rho^s} \int_{S^{n-1}} P_j(u) dS(u) = \frac{1}{j+n-s} \int_{S^{n-1}} P_j(u) dS(u).$$

Lemma 7.3 now gives the result. □

7.2 Holomorphy of certain series

Before stating the main result of this section, we give first in the following some preliminaries from Diophantine approximation theory:

Definition 7.4. (i) Let $\delta > 0$. A vector $a \in \mathbb{R}^n$ is said to be δ -badly approximable if there exists $c > 0$ such that $|q \cdot a - m| \geq c |q|^{-\delta}, \forall q \in \mathbb{Z}^n \setminus \{0\}$ and $\forall m \in \mathbb{Z}$.

We note $\mathcal{BV}(\delta)$ the set of δ -badly approximable vectors and $\mathcal{BV} := \cup_{\delta > 0} \mathcal{BV}(\delta)$ the set of badly approximable vectors.

(ii) A matrix $\Theta \in \mathcal{M}_n(\mathbb{R})$ (real $n \times n$ matrices) will be said to be badly approximable if there exists $u \in \mathbb{Z}^n$ such that ${}^t\Theta(u)$ is a badly approximable vector of \mathbb{R}^n .

Remark. A classical result from Diophantine approximation asserts that for $\delta > n$, the Lebesgue measure of $\mathbb{R}^n \setminus \mathcal{BV}(\delta)$ is zero (i.e almost any element of \mathbb{R}^n is δ -badly approximable.)

Let $\Theta \in \mathcal{M}_n(\mathbb{R})$. If its row of index i is a badly approximable vector of \mathbb{R}^n (i.e. if $L_i \in \mathcal{BV}$) then ${}^t\Theta(e_i) \in \mathcal{BV}$ and thus Θ is a badly approximable matrix. It follows that almost any matrix of $\mathcal{M}_n(\mathbb{R}) \approx \mathbb{R}^{n^2}$ is badly approximable.

The goal of this section is to show the following

Theorem 7.5. Let $P \in \mathbb{C}[X_1, \dots, X_n]$ be a homogeneous polynomial of degree d and let b be in $\mathcal{S}(\mathbb{Z}^n \times \dots \times \mathbb{Z}^n)$ (q times, $q \in \mathbb{N}$). Then,

- (i) Let $a \in \mathbb{R}^n$. We define $f_a(s) := \sum_{k \in \mathbb{Z}^n} \frac{P(k)}{|k|^s} e^{2\pi i k \cdot a}$.

1. If $a \in \mathbb{Z}^n$, then f_a has a meromorphic continuation to the whole complex plane \mathbb{C} . Moreover if S is the unit sphere and dS its Lebesgue measure, then f_a is not entire if and only if $\int_{u \in S^{n-1}} P(u) dS(u) \neq 0$. In that case, f_a has only a simple pole at the point $d+n$, with $\text{Res}_{s=d+n} f_a(s) = \int_{u \in S^{n-1}} P(u) dS(u)$.

2. If $a \in \mathbb{R}^n \setminus \mathbb{Z}^n$, then $f_a(s)$ extends holomorphically to the whole complex plane \mathbb{C} .

(ii) Suppose that $\Theta \in \mathcal{M}_n(\mathbb{R})$ is badly approximable. For any $(\varepsilon_i)_i \in \{-1, 0, 1\}^q$, the function

$$g(s) := \sum_{l \in (\mathbb{Z}^n)^q} b(l) f_{\Theta \sum_i \varepsilon_i l_i}(s)$$

extends meromorphically to the whole complex plane \mathbb{C} with only one possible pole on $s = d+n$.

Moreover, if we set $\mathcal{Z} := \{l \in (\mathbb{Z}^n)^q \mid \sum_{i=1}^q \varepsilon_i l_i = 0\}$ and $V := \sum_{l \in \mathcal{Z}} b(l)$, then

1. If $V \int_{S^{n-1}} P(u) dS(u) \neq 0$, then $s = d+n$ is a simple pole of $g(s)$ and

$$\text{Res}_{s=d+n} g(s) = V \int_{u \in S^{n-1}} P(u) dS(u).$$

2. If $V \int_{S^{n-1}} P(u) dS(u) = 0$, then $g(s)$ extends holomorphically to the whole complex plane \mathbb{C} .

(iii) Suppose that $\Theta \in \mathcal{M}_n(\mathbb{R})$ is badly approximable. For any $(\varepsilon_i)_i \in \{-1, 0, 1\}^q$, the function

$$g_0(s) := \sum_{l \in (\mathbb{Z}^n)^q \setminus \mathcal{Z}} b(l) f_{\Theta \sum_{i=1}^q \varepsilon_i l_i}(s)$$

where $\mathcal{Z} := \{l \in (\mathbb{Z}^n)^q \mid \sum_{i=1}^q \varepsilon_i l_i = 0\}$ extends holomorphically to the whole complex plane \mathbb{C} .

Proof of Theorem 7.5:

First we remark that

If $a \in \mathbb{Z}^n$ then $f_a(s) = \sum'_{k \in \mathbb{Z}^n} \frac{P(k)}{|k|^s}$. So, the point (i.1) follows from Theorem 7.1;

$g(s) := \sum_{l \in (\mathbb{Z}^n)^q \setminus \mathcal{Z}} b(l) f_{\Theta \sum_i \varepsilon_i l_i}(s) + (\sum_{l \in \mathcal{Z}} b(l)) \sum'_{k \in \mathbb{Z}^n} \frac{P(k)}{|k|^s}$. Thus, the point (ii) rises easily from (iii) and Theorem 7.1.

So, to complete the proof, it remains to prove the items (i.2) and (iii).

The direct proof of (i.2) is easy but is not sufficient to deduce (iii) of which the proof is more delicate and requires a more precise (i.e. more effective) version of (i.2). The next lemma gives such crucial version, but before, let us give some notations:

$$\mathcal{F} := \left\{ \frac{P(X)}{(X_1^2 + \dots + X_n^2 + 1)^{r/2}} \mid P(X) \in \mathbb{C}[X_1, \dots, X_n] \text{ and } r \in \mathbb{N}_0 \right\}.$$

We set $g = \deg(G) = \deg(P) - r \in \mathbb{Z}$, the degree of $G = \frac{P(X)}{(X_1^2 + \dots + X_n^2 + 1)^{r/2}} \in \mathcal{F}$.

By convention, we set $\deg(0) = -\infty$.

Lemma 7.6. *Let $a \in \mathbb{R}^n$. We assume that $d(a.u, \mathbb{Z}) := \inf_{m \in \mathbb{Z}} |a.u - m| > 0$ for some $u \in \mathbb{Z}^n$. For all $G \in \mathcal{F}$, we define formally,*

$$F_0(G; a; s) := \sum'_{k \in \mathbb{Z}^n} \frac{G(k)}{|k|^s} e^{2\pi i k.a} \quad \text{and} \quad F_1(G; a; s) := \sum_{k \in \mathbb{Z}^n} \frac{G(k)}{(|k|^2 + 1)^{s/2}} e^{2\pi i k.a}.$$

Then for all $N \in \mathbb{N}$, $G \in \mathcal{F}$ and $i \in \{0, 1\}$, there exist positive constants $C_i := C_i(G, N, u)$, $B_i := B_i(G, N, u)$ and $A_i := A_i(G, N, u)$ such that $s \mapsto F_i(G; a; s)$ extends holomorphically to the half-plane $\{\Re(s) > -N\}$ and verifies in it:

$$F_i(G; a; s) \leq C_i (1 + |s|)^{B_i} \left(d(a.u, \mathbb{Z}) \right)^{-A_i}.$$

Remark 7.7. *The important point here is that we obtain an explicit bound of $F_i(G; \alpha; s)$ in $\{\Re(s) > -N\}$ which depends on the vector a only through $d(a, u, \mathbb{Z})$, so depends on u and indirectly on a (in the sequel, a will vary.) In particular the constants $C_i := C_i(G, N, u)$, $B_i = B_i(G, N)$ and $A_i := A_i(G, N)$ do not depend on the vector a but only on u . This is crucial for the proof of items (ii) and (iii) of Theorem 7.5!*

7.2.1 Proof of Lemma 7.6 for $i = 1$:

Let $N \in \mathbb{N}_0$ be a fixed integer, and set $g_0 := n + N + 1$.

We will prove Lemma 7.6 by induction on $g = \deg(G) \in \mathbb{Z}$. More precisely, in order to prove case $i = 1$, it suffices to prove that:

Lemma 7.6 is true for all $G \in \mathcal{F}$ with $\deg(G) \leq -g_0$.

Let $g \in \mathbb{Z}$ with $g \geq -g_0 + 1$. If Lemma 7.6 is true for all $G \in \mathcal{F}$ such that $\deg(G) \leq g - 1$,

then it is also true for all $G \in \mathcal{F}$ satisfying $\deg(G) = g$.

- Step 1: Checking Lemma 7.6 for $\deg(G) \leq -g_0 := -(n + N + 1)$.

Let $G(X) = \frac{P(X)}{(X_1^2 + \dots + X_n^2 + 1)^{r/2}} \in \mathcal{F}$ with $\deg(G) \leq -g_0$. It is easy to see that we have uniformly in $s = \sigma + i\tau \in \mathbb{C}$ and in $k \in \mathbb{Z}^n$:

$$\frac{|G(k) e^{2\pi i k \cdot a}|}{(|k|^2 + 1)^{\sigma/2}} = \frac{|P(k)|}{(|k|^2 + 1)^{(r+\sigma)/2}} \ll_G \frac{1}{(|k|^2 + 1)^{(r+\sigma - \deg(P))/2}} \ll_G \frac{1}{(|k|^2 + 1)^{(\sigma - \deg(G))/2}} \ll_G \frac{1}{(|k|^2 + 1)^{(\sigma + g_0)/2}}.$$

It follows that $F_1(G; a; s) = \sum_{k \in \mathbb{Z}^n} \frac{G(k)}{(|k|^2 + 1)^{s/2}} e^{2\pi i k \cdot a}$ converges absolutely and defines a holomorphic function in the half plane $\{\sigma > -N\}$. Therefore, we have for any $s \in \{\Re(s) > -N\}$:

$$|F_1(G; a; s)| \ll_G \sum_{k \in \mathbb{Z}^n} \frac{1}{(|k|^2 + 1)^{(-N + g_0)/2}} \ll_G \sum_{k \in \mathbb{Z}^n} \frac{1}{(|k|^2 + 1)^{(n+1)/2}} \ll_G 1.$$

Thus, Lemma 7.6 is true when $\deg(G) \leq -g_0$.

- Step 2: Induction.

Now let $g \in \mathbb{Z}$ satisfying $g \geq -g_0 + 1$ and suppose that Lemma 7.6 is valid for all $G \in \mathcal{F}$ with $\deg(G) \leq g - 1$. Let $G \in \mathcal{F}$ with $\deg(G) = g$. We will prove that G also verifies conclusions of Lemma 7.6:

There exist $P \in \mathbb{C}[X_1, \dots, X_n]$ of degree $d \geq 0$ and $r \in \mathbb{N}_0$ such that $G(X) = \frac{P(X)}{(X_1^2 + \dots + X_n^2 + 1)^{r/2}}$ and $g = \deg(G) = d - r$.

Since $G(k) \ll (|k|^2 + 1)^{g/2}$ uniformly in $k \in \mathbb{Z}^n$, we deduce that $F_1(G; a; s)$ converges absolutely in $\{\sigma = \Re(s) > n + g\}$.

Since $k \mapsto k + u$ is a bijection from \mathbb{Z}^n into \mathbb{Z}^n , it follows that we also have for $\Re(s) > n + g$

$$\begin{aligned} F_1(G; a; s) &= \sum_{k \in \mathbb{Z}^n} \frac{P(k)}{(|k|^2 + 1)^{(s+r)/2}} e^{2\pi i k \cdot a} = \sum_{k \in \mathbb{Z}^n} \frac{P(k+u)}{(|k+u|^2 + 1)^{(s+r)/2}} e^{2\pi i (k+u) \cdot a} \\ &= e^{2\pi i u \cdot a} \sum_{k \in \mathbb{Z}^n} \frac{P(k+u)}{(|k|^2 + 2k \cdot u + |u|^2 + 1)^{(s+r)/2}} e^{2\pi i k \cdot a} \\ &= e^{2\pi i u \cdot a} \sum_{\alpha \in \mathbb{N}_0^n; |\alpha|_1 = \alpha_1 + \dots + \alpha_n \leq d} \frac{u^\alpha}{\alpha!} \sum_{k \in \mathbb{Z}^n} \frac{\partial^\alpha P(k)}{(|k|^2 + 2k \cdot u + |u|^2 + 1)^{(s+r)/2}} e^{2\pi i k \cdot a} \\ &= e^{2\pi i u \cdot a} \sum_{|\alpha|_1 \leq d} \frac{u^\alpha}{\alpha!} \sum_{k \in \mathbb{Z}^n} \frac{\partial^\alpha P(k)}{(|k|^2 + 1)^{(s+r)/2}} \left(1 + \frac{2k \cdot u + |u|^2}{(|k|^2 + 1)}\right)^{-(s+r)/2} e^{2\pi i k \cdot a}. \end{aligned}$$

Let $M := \sup(N + n + g, 0) \in \mathbb{N}_0$. We have uniformly in $k \in \mathbb{Z}^n$

$$\left(1 + \frac{2k \cdot u + |u|^2}{(|k|^2 + 1)}\right)^{-(s+r)/2} = \sum_{j=0}^M \binom{-(s+r)/2}{j} \frac{(2k \cdot u + |u|^2)^j}{(|k|^2 + 1)^j} + O_{M,u} \left(\frac{(1+|s|)^{M+1}}{(|k|^2 + 1)^{(M+1)/2}} \right).$$

Thus, for $\sigma = \Re(s) > n + d$,

$$\begin{aligned} F_1(G; a; s) &= e^{2\pi i u \cdot a} \sum_{|\alpha|_1 \leq d} \frac{u^\alpha}{\alpha!} \sum_{k \in \mathbb{Z}^n} \frac{\partial^\alpha P(k)}{(|k|^2 + 1)^{(s+r)/2}} \left(1 + \frac{2k \cdot u + |u|^2}{(|k|^2 + 1)}\right)^{-(s+r)/2} e^{2\pi i k \cdot a} \\ &= e^{2\pi i u \cdot a} \sum_{|\alpha|_1 \leq d} \sum_{j=0}^M \frac{u^\alpha}{\alpha!} \binom{-(s+r)/2}{j} \sum_{k \in \mathbb{Z}^n} \frac{\partial^\alpha P(k) (2k \cdot u + |u|^2)^j}{(|k|^2 + 1)^{(s+r+2j)/2}} e^{2\pi i k \cdot a} \\ &\quad + O_{G,M,u} \left((1+|s|)^{M+1} \sum_{k \in \mathbb{Z}^n} \frac{1}{(|k|^2 + 1)^{(\sigma+M+1-g)/2}} \right). \end{aligned} \quad (80)$$

Set $I := \{(\alpha, j) \in \mathbb{N}_0^n \times \{0, \dots, M\} \mid |\alpha|_1 \leq d\}$ and $I^* := I \setminus \{(0, 0)\}$.

Set also $G_{(\alpha,j);u}(X) := \frac{\partial^\alpha P(X) (2X \cdot u + |u|^2)^j}{(|X|^2 + 1)^{(r+2j)/2}} \in \mathcal{F}$ for all $(\alpha, j) \in I^*$.

Since $M \geq N + n + g$, it follows from (80) that

$$(1 - e^{2\pi i u \cdot a}) F_1(G; a; s) = e^{2\pi i u \cdot a} \sum_{(\alpha,j) \in I^*} \frac{u^\alpha}{\alpha!} \binom{-(s+r)/2}{j} F_1(G_{(\alpha,j);u}; \alpha; s) + R_N(G; a; u; s) \quad (81)$$

where $s \mapsto R_N(G; a; u; s)$ is a holomorphic function in the half plane $\{\sigma = \Re(s) > -N\}$, in which it satisfies the bound $R_N(G; a; u; s) \ll_{G,N,u} 1$.

Moreover it is easy to see that, for any $(\alpha, j) \in I^*$,

$$\deg(G_{(\alpha,j);u}) = \deg(\partial^\alpha P) + j - (r + 2j) \leq d - |\alpha|_1 + j - (r + 2j) = g - |\alpha|_1 - j \leq g - 1.$$

Relation (81) and the induction hypothesis imply then that

$$(1 - e^{2\pi i u \cdot a}) F_1(G; a; s) \text{ verifies the conclusions of Lemma 7.6.} \quad (82)$$

Since $|1 - e^{2\pi i u \cdot a}| = 2|\sin(\pi u \cdot a)| \geq d(u \cdot a, \mathbb{Z})$, then (82) implies that $F_1(G; a; s)$ satisfies conclusions of Lemma 7.6. This completes the induction and the proof for $i = 1$.

7.2.2 Proof of Lemma 7.6 for $i = 0$:

Let $N \in \mathbb{N}$ be a fixed integer. Let $G(X) = \frac{P(X)}{(X_1^2 + \dots + X_n^2 + 1)^{r/2}} \in \mathcal{F}$ and $g = \deg(G) = d - r$ where $d \geq 0$ is the degree of the polynomial P . Set also $M := \sup(N + g + n, 0) \in \mathbb{N}_0$.

Since $P(k) \ll |k|^d$ for $k \in \mathbb{Z}^n \setminus \{0\}$, it follows that $F_0(G; a; s)$ and $F_1(G; a; s)$ converge absolutely in the half plane $\{\sigma = \Re(s) > n + g\}$.

Moreover, we have for $s = \sigma + i\tau \in \mathbb{C}$ with $\sigma > n + g$:

$$\begin{aligned}
F_0(G; a; s) &= \sum_{k \in \mathbb{Z}^n \setminus \{0\}} \frac{G(k)}{(|k|^2 + 1)^{s/2}} e^{2\pi i k \cdot a} = \sum'_{k \in \mathbb{Z}^n} \frac{G(k)}{(|k|^2 + 1)^{s/2}} \left(1 - \frac{1}{|k|^2 + 1}\right)^{-s/2} e^{2\pi i k \cdot a} \\
&= \sum'_{k \in \mathbb{Z}^n} \sum_{j=0}^M \binom{-s/2}{j} (-1)^j \frac{G(k)}{(|k|^2 + 1)^{(s+2j)/2}} e^{2\pi i k \cdot a} \\
&\quad + O_M\left((1 + |s|)^{M+1} \sum'_{k \in \mathbb{Z}^n} \frac{|G(k)|}{(|k|^2 + 1)^{(\sigma+2M+2)/2}}\right) \\
&= \sum_{j=0}^M \binom{-s/2}{j} (-1)^j F_1(G; a; s + 2j) \\
&\quad + O_M\left[(1 + |s|)^{M+1} \left(1 + \sum'_{k \in \mathbb{Z}^n} \frac{|G(k)|}{(|k|^2 + 1)^{(\sigma+2M+2)/2}}\right)\right]. \tag{83}
\end{aligned}$$

In addition we have uniformly in $s = \sigma + i\tau \in \mathbb{C}$ with $\sigma > -N$,

$$\sum'_{k \in \mathbb{Z}^n} \frac{|G(k)|}{(|k|^2 + 1)^{(\sigma+2M+2)/2}} \ll \sum'_{k \in \mathbb{Z}^n} \frac{|k|^g}{(|k|^2 + 1)^{(-N+2M+2)/2}} \ll \sum'_{k \in \mathbb{Z}^n} \frac{1}{|k|^{n+1}} < +\infty.$$

So (83) and Lemma 7.6 for $i = 1$ imply that Lemma 7.6 is also true for $i = 0$. This completes the proof of Lemma 7.6.

7.2.3 Proof of item (i.2) of Theorem 7.5:

Since $a \in \mathbb{R}^n \setminus \mathbb{Z}^n$, there exists $i_0 \in \{1, \dots, n\}$ with $a_{i_0} \notin \mathbb{Z}$. So $d(a \cdot e_{i_0}, \mathbb{Z}) = d(a_{i_0}, \mathbb{Z}) > 0$. Therefore, a satisfies the assumption of Lemma 7.6 with $u = e_{i_0}$. Thus, for all $N \in \mathbb{N}$, $s \mapsto f_a(s) = F_0(P; a; s)$ has a holomorphic continuation to the half-plane $\{\Re(s) > -N\}$. It follows, by letting $N \rightarrow \infty$, that $s \mapsto f_a(s)$ has a holomorphic continuation to the whole complex plane \mathbb{C} .

7.2.4 Proof of item (iii) of Theorem 7.5:

Let $\Theta \in \mathcal{M}_n(\mathbb{R})$, $(\varepsilon_i)_i \in \{-1, 0, 1\}^q$ and $b \in \mathcal{S}(\mathbb{Z}^n \times \mathbb{Z}^n)$. We assume that Θ is a badly approximable matrix. Set $\mathcal{Z} := \{l = (l_1, \dots, l_q) \in (\mathbb{Z}^n)^q \mid \sum_i \varepsilon_i l_i = 0\}$ and $P \in \mathbb{C}[X_1, \dots, X_n]$ of degree $d \geq 0$.

It is easy to see that for $\sigma > n + d$:

$$\begin{aligned}
\sum_{l \in (\mathbb{Z}^n)^q \setminus \mathcal{Z}} |b(l)| \sum'_{k \in \mathbb{Z}^n} \frac{|P(k)|}{|k|^\sigma} |e^{2\pi i k \cdot \Theta \sum_i \varepsilon_i l_i}| &\ll_P \sum_{l \in (\mathbb{Z}^n)^q \setminus \mathcal{Z}} |b(l)| \sum'_{k \in \mathbb{Z}^n} \frac{1}{|k|^{\sigma-d}} \ll_{P, \sigma} \sum_{l \in (\mathbb{Z}^n)^q \setminus \mathcal{Z}} |b(l)| \\
&< +\infty.
\end{aligned}$$

So

$$g_0(s) := \sum_{l \in (\mathbb{Z}^n)^q \setminus \mathcal{Z}} b(l) f_{\Theta \sum_i \varepsilon_i l_i}(s) = \sum_{l \in (\mathbb{Z}^n)^q \setminus \mathcal{Z}} b(l) \sum'_{k \in \mathbb{Z}^n} \frac{P(k)}{|k|^s} e^{2\pi i k \cdot \Theta \sum_i \varepsilon_i l_i}$$

converges absolutely in the half plane $\{\Re(s) > n + d\}$.

Moreover with the notations of Lemma 7.6, we have for all $s = \sigma + i\tau \in \mathbb{C}$ with $\sigma > n + d$:

$$g_0(s) = \sum_{l \in (\mathbb{Z}^n)^q \setminus \mathcal{Z}} b(l) f_{\Theta \sum_i \varepsilon_i l_i}(s) = \sum_{l \in (\mathbb{Z}^n)^q \setminus \mathcal{Z}} b(l) F_0(P; \Theta \sum_i \varepsilon_i l_i; s) \tag{84}$$

But Θ is badly approximable, so there exists $u \in \mathbb{Z}^n$ and $\delta, c > 0$ such

$$|q \cdot {}^t\Theta u - m| \geq c(1 + |q|)^{-\delta}, \forall q \in \mathbb{Z}^n \setminus \{0\}, \forall m \in \mathbb{Z}.$$

We deduce that $\forall l \in (\mathbb{Z}^n)^q \setminus \mathcal{Z}$,

$$|(\Theta \sum_i \varepsilon_i l_i) \cdot u - m| = |(\sum_i \varepsilon_i l_i) \cdot {}^t\Theta u - m| \geq c(1 + |\sum_i \varepsilon_i l_i|)^{-\delta} \geq c(1 + |l|)^{-\delta}.$$

It follows that there exists $u \in \mathbb{Z}^n$, $\delta > 0$ and $c > 0$ such that

$$\forall l \in (\mathbb{Z}^n)^q \setminus \mathcal{Z}, \quad d((\Theta \sum_i \varepsilon_i l_i) \cdot u; \mathbb{Z}) \geq c(1 + |l|)^{-\delta}. \quad (85)$$

Therefore, for any $l \in (\mathbb{Z}^n)^q \setminus \mathcal{Z}$, the vector $a = \Theta \sum_i \varepsilon_i l_i$ verifies the assumption of Lemma 7.6 with the same u . Moreover δ and c in (85) are also independent on l .

We fix now $N \in \mathbb{N}$. Lemma 7.6 implies that there exist positive constants $C_0 := C_0(P, N, u)$, $B_0 := B_0(P, N, u)$ and $A_0 := A_0(P, N, u)$ such that for all $l \in (\mathbb{Z}^n)^q \setminus \mathcal{Z}$, $s \mapsto F_0(P; \Theta \sum_i \varepsilon_i l_i; s)$ extends holomorphically to the half plane $\{\Re(s) > -N\}$ and verifies in it the bound

$$F_0(P; \Theta \sum_i \varepsilon_i l_i; s) \leq C_0(1 + |s|)^{B_0} d((\Theta \sum_i \varepsilon_i l_i) \cdot u; \mathbb{Z})^{-A_0}.$$

This and (85) imply that for any compact set K included in the half plane $\{\Re(s) > -N\}$, there exist two constants $C := C(P, N, c, \delta, u, K)$ and $D := D(P, N, c, \delta, u)$ (independent on $l \in (\mathbb{Z}^n)^q \setminus \mathcal{Z}$) such that

$$\forall s \in K \text{ and } \forall l \in (\mathbb{Z}^n)^q \setminus \mathcal{Z}, \quad F_0(P; \Theta \sum_i \varepsilon_i l_i; s) \leq C(1 + |l|)^D. \quad (86)$$

It follows that $s \mapsto \sum_{l \in (\mathbb{Z}^n)^q \setminus \mathcal{Z}} b(l) F_0(P; \Theta \sum_i \varepsilon_i l_i; s)$ has a holomorphic continuation to the half plane $\{\Re(s) > -N\}$.

This and (84) imply that $s \mapsto g_0(s) = \sum_{l \in (\mathbb{Z}^n)^q \setminus \mathcal{Z}} b(l) f_{\Theta \sum_i \varepsilon_i l_i}(s)$ has a holomorphic continuation to $\{\Re(s) > -N\}$. Since N is an arbitrary integer, by letting $N \rightarrow \infty$, it follows that $s \mapsto g_0(s)$ has a holomorphic continuation to the whole complex plane \mathbb{C} which completes the proof of the theorem.

Remark 7.8. By equation (82), we see that a Diophantine condition is sufficient to get Lemma 7.6. Our Diophantine condition appears also (in equivalent form) in Connes [23, Prop. 49] (see Remark 4.2 below). The following heuristic argument shows that our condition seems to be necessary in order to get the result of Theorem 7.5:

For simplicity we assume $n = 1$ (but the argument extends easily to any n).

Let $\theta \in \mathbb{R} \setminus \mathbb{Q}$. We know that for any $l \in \mathbb{Z} \setminus \{0\}$,

$$g_{\theta l}(s) := \sum'_{k \in \mathbb{Z}} \frac{e^{2\pi i \theta l k}}{|k|^s} = \frac{\pi^{s-1/2}}{\Gamma(\frac{1-s}{2})} \Gamma(\frac{s}{2}) h_{\theta l}(1-s) \text{ where } h_{\theta l}(s) := \sum'_{k \in \mathbb{Z}} \frac{1}{|\theta l + k|^s}.$$

So, for any $(a_l) \in \mathcal{S}(\mathbb{Z})$, the existence of meromorphic continuation of $g_0(s) := \sum'_{l \in \mathbb{Z}} a_l g_{\theta l}(s)$ is equivalent to the existence of meromorphic continuation of

$$h_0(s) := \sum'_{l \in \mathbb{Z}} a_l h_{\theta l}(s) = \sum'_{l \in \mathbb{Z}} a_l \sum'_{k \in \mathbb{Z}} \frac{1}{|\theta l + k|^s}.$$

So, for at least one $\sigma_0 \in \mathbb{R}$, we must have $\frac{|a_l|}{|\theta l + k|^{\sigma_0}} = O(1)$ uniformly in $k, l \in \mathbb{Z}^*$.

It follows that for any $(a_l) \in \mathcal{S}(\mathbb{Z})$, $|\theta l + k| \gg |a_l|^{1/\sigma_0}$ uniformly in $k, l \in \mathbb{Z}^*$. Therefore, our Diophantine condition seems to be necessary.

7.2.5 Commutation between sum and residue

Let $p \in \mathbb{N}$. Recall that $\mathcal{S}((\mathbb{Z}^n)^p)$ is the set of the Schwartz sequences on $(\mathbb{Z}^n)^p$. In other words, $b \in \mathcal{S}((\mathbb{Z}^n)^p)$ if and only if for all $r \in \mathbb{N}_0$, $(1 + |l_1|^2 + \dots + |l_p|^2)^r |b(l_1, \dots, l_p)|^2$ is bounded on $(\mathbb{Z}^n)^p$. We note that if $Q \in \mathbb{R}[X_1, \dots, X_{np}]$ is a polynomial, $(a_j) \in \mathcal{S}(\mathbb{Z}^n)^p$, $b \in \mathcal{S}(\mathbb{Z}^n)$ and ϕ a real-valued function, then $l := (l_1, \dots, l_p) \mapsto \tilde{a}(l) b(-\widehat{l}_p) Q(l) e^{i\phi(l)}$ is a Schwartz sequence on $(\mathbb{Z}^n)^p$, where

$$\begin{aligned}\tilde{a}(l) &:= a_1(l_1) \cdots a_p(l_p), \\ \widehat{l}_i &:= l_1 + \dots + l_i.\end{aligned}$$

In the following, we will use several times the fact that for any $(k, l) \in (\mathbb{Z}^n)^2$ such that $k \neq 0$ and $k \neq -l$, we have

$$\frac{1}{|k+l|^2} = \frac{1}{|k|^2} - \frac{2k \cdot l + |l|^2}{|k|^2 |k+l|^2}. \quad (87)$$

Lemma 7.9. *There exists a polynomial $P \in \mathbb{R}[X_1, \dots, X_p]$ of degree $4p$ and with positive coefficients such that for any $k \in \mathbb{Z}^n$, and $l := (l_1, \dots, l_p) \in (\mathbb{Z}^n)^p$ such that $k \neq 0$ and $k \neq -\widehat{l}_i$ for all $1 \leq i \leq p$, the following holds:*

$$\frac{1}{|k + \widehat{l}_1|^2 \dots |k + \widehat{l}_p|^2} \leq \frac{1}{|k|^{2p}} P(|l_1|, \dots, |l_p|).$$

Proof. Let's fix i such that $1 \leq i \leq p$. Using two times (87), Cauchy–Schwarz inequality and the fact that $|k + \widehat{l}_i|^2 \geq 1$, we get

$$\begin{aligned}\frac{1}{|k + \widehat{l}_i|^2} &\leq \frac{1}{|k|^2} + \frac{2|k||\widehat{l}_i| + |\widehat{l}_i|^2}{|k|^4} + \frac{(2|k||\widehat{l}_i| + |\widehat{l}_i|^2)^2}{|k|^4 |k + \widehat{l}_i|^2} \\ &\leq \frac{1}{|k|^2} + \frac{2}{|k|^3} |\widehat{l}_i| + \left(\frac{1}{|k|^4} + \frac{4}{|k|^2} \right) |\widehat{l}_i|^2 + \frac{4}{|k|^3} |\widehat{l}_i|^3 + \frac{1}{|k|^4} |\widehat{l}_i|^4.\end{aligned}$$

Since $|k| \geq 1$, and $|\widehat{l}_i|^j \leq |\widehat{l}_i|^4$ if $1 \leq j \leq 4$, we find

$$\begin{aligned}\frac{1}{|k + \widehat{l}_i|^2} &\leq \frac{5}{|k|^2} \sum_{j=0}^4 |\widehat{l}_i|^j \leq \frac{5}{|k|^2} (1 + 4|\widehat{l}_i|^4) \leq \frac{5}{|k|^2} (1 + 4(\sum_{j=1}^p |l_j|^4)^2), \\ \frac{1}{|k + \widehat{l}_1|^2 \dots |k + \widehat{l}_p|^2} &\leq \frac{5^p}{|k|^{2p}} (1 + 4(\sum_{j=1}^p |l_j|^4)^2)^p.\end{aligned}$$

Taking $P(X_1, \dots, X_p) := 5^p (1 + 4(\sum_{j=1}^p X_j^4)^2)^p$ now gives the result. \square

Lemma 7.10. *Let $b \in \mathcal{S}((\mathbb{Z}^n)^p)$, $p \in \mathbb{N}$, $P_j \in \mathbb{R}[X_1, \dots, X_n]$ be a homogeneous polynomial function of degree j , $k \in \mathbb{Z}^n$, $l := (l_1, \dots, l_p) \in (\mathbb{Z}^n)^p$, $r \in \mathbb{N}_0$, ϕ be a real-valued function on $\mathbb{Z}^n \times (\mathbb{Z}^n)^p$ and*

$$h(s, k, l) := \frac{b(l) P_j(k) e^{i\phi(k, l)}}{|k|^{s+r} |k + \widehat{l}_1|^2 \dots |k + \widehat{l}_p|^2},$$

with $h(s, k, l) := 0$ if, for $k \neq 0$, one of the denominators is zero.

For all $s \in \mathbb{C}$ such that $\Re(s) > n + j - r - 2p$, the series

$$H(s) := \sum'_{(k, l) \in (\mathbb{Z}^n)^{p+1}} h(s, k, l)$$

is absolutely summable. In particular,

$$\sum'_{k \in \mathbb{Z}^n} \sum_{l \in (\mathbb{Z}^n)^p} h(s, k, l) = \sum_{l \in (\mathbb{Z}^n)^p} \sum'_{k \in \mathbb{Z}^n} h(s, k, l).$$

Proof. Let $s = \sigma + i\tau \in \mathbb{C}$ such that $\sigma > n + j - r - 2p$. By Lemma 7.9 we get, for $k \neq 0$,

$$|h(s, k, l)| \leq |b(l) P_j(k)| |k|^{-r-\sigma-2p} P(l),$$

where $P(l) := P(|l_1|, \dots, |l_p|)$ and P is a polynomial of degree $4p$ with positive coefficients. Thus, $|h(s, k, l)| \leq F(l) G(k)$ where $F(l) := |b(l)| P(l)$ and $G(k) := |P_j(k)| |k|^{-r-\sigma-2p}$. The summability of $\sum_{l \in (\mathbb{Z}^n)^p} F(l)$ is implied by the fact that $b \in \mathcal{S}((\mathbb{Z}^n)^p)$. The summability of $\sum'_{k \in \mathbb{Z}^n} G(k)$ is a consequence of the fact that $\sigma > n + j - r - 2p$. Finally, as a product of two summable series, $\sum_{k,l} F(l) G(k)$ is a summable series, which proves that $\sum_{k,l} h(s, k, l)$ is also absolutely summable. \square

Definition 7.11. Let f be a function on $D \times (\mathbb{Z}^n)^p$ where D is an open neighborhood of 0 in \mathbb{C} .

We say that f satisfies (H1) if and only if there exists $\rho > 0$ such that

(i) for any l , $s \mapsto f(s, l)$ extends as a holomorphic function on U_ρ , where U_ρ is the open disk of center 0 and radius ρ ,

(ii) if $\|H(\cdot, l)\|_{\infty, \rho} := \sup_{s \in U_\rho} |H(s, l)|$, the series $\sum_{l \in (\mathbb{Z}^n)^p} \|H(\cdot, l)\|_{\infty, \rho}$ is summable.

We say that f satisfies (H2) if and only if there exists $\rho > 0$ such that

(i) for any l , $s \mapsto f(s, l)$ extends as a holomorphic function on $U_\rho - \{0\}$,

(ii) for any δ such that $0 < \delta < \rho$, the series $\sum_{l \in (\mathbb{Z}^n)^p} \|H(\cdot, l)\|_{\infty, \delta, \rho}$ is summable, where $\|H(\cdot, l)\|_{\infty, \delta, \rho} := \sup_{\delta < |s| < \rho} |H(s, l)|$.

Remark 7.12. Note that (H1) implies (H2). Moreover, if f satisfies (H1) (resp. (H2)) for $\rho > 0$, then it is straightforward to check that $f : s \mapsto \sum_{l \in (\mathbb{Z}^n)^p} f(s, l)$ extends as an holomorphic function on U_ρ (resp. on $U_\rho \setminus \{0\}$).

Corollary 7.13. With the same notations of Lemma 7.10, suppose that $r + 2p - j > n$, then, the function $H(s, l) := \sum'_{k \in \mathbb{Z}^n} h(s, k, l)$ satisfies (H1).

Proof. (i) Let's fix $\rho > 0$ such that $\rho < r + 2p - j - n$. Since $r + 2p - j > n$, U_ρ is inside the half-plane of absolute convergence of the series defined by $H(s, l)$. Thus, $s \mapsto H(s, l)$ is holomorphic on U_ρ .

(ii) Since $|k|^{-s} \leq |k|^\rho$ for all $s \in U_\rho$ and $k \in \mathbb{Z}^n \setminus \{0\}$, we get as in the above proof

$$|h(s, k, l)| \leq |b(l) P_j(k)| |k|^{-r+\rho-2p} P(|l_1|, \dots, |l_p|).$$

Since $\rho < r + 2p - j - n$, the series $\sum'_{k \in \mathbb{Z}^n} |P_j(k)| |k|^{-r+\rho-2p}$ is summable.

Thus, $\|H(\cdot, l)\|_{\infty, \rho} \leq K F(l)$ where $K := \sum_k' |P_j(k)| |k|^{-r+\rho-2p} < \infty$. We have already seen that the series $\sum_l F(l)$ is summable, so we get the result. \square

We note that if f and g both satisfy (H1) (or (H2)), then so does $f + g$. In the following, we will use the equivalence relation

$$f \sim g \iff f - g \text{ satisfies (H1)}.$$

Lemma 7.14. Let f and g be two functions on $D \times (\mathbb{Z}^n)^p$ where D is an open neighborhood of 0 in \mathbb{C} , such that $f \sim g$ and such that g satisfies (H2). Then

$$\operatorname{Res}_{s=0} \sum_{l \in (\mathbb{Z}^n)^p} f(s, l) = \sum_{l \in (\mathbb{Z}^n)^p} \operatorname{Res}_{s=0} g(s, l).$$

Proof. Since $f \sim g$, f satisfies (H2) for a certain $\rho > 0$. Let's fix η such that $0 < \eta < \rho$ and define C_η as the circle of center 0 and radius η . We have

$$\operatorname{Res}_{s=0} g(s, l) = \operatorname{Res}_{s=0} f(s, l) = \frac{1}{2\pi i} \oint_{C_\eta} f(s, l) ds = \int_I u(t, l) dt.$$

where $I = [0, 2\pi]$ and $u(t, l) := \frac{1}{2\pi} \eta e^{it} f(\eta e^{it}, l)$. The fact that f satisfies (H2) entails that the series $\sum_{l \in (\mathbb{Z}^n)^p} \|f(\cdot, l)\|_{\infty, C_\eta}$ is summable. Thus, since $\|u(\cdot, l)\|_\infty = \frac{1}{2\pi} \eta \|f(\cdot, l)\|_{\infty, C_\eta}$, the series $\sum_{l \in (\mathbb{Z}^n)^p} \|u(\cdot, l)\|_\infty$ is summable, so, $\int_I \sum_{l \in (\mathbb{Z}^n)^p} u(t, l) dt = \sum_{l \in (\mathbb{Z}^n)^p} \int_I u(t, l) dt$ which gives the result. \square

7.3 Computation of residues of zeta functions

Since, we will have to compute residues of series, let us introduce the following

Definition 7.15.

$$\begin{aligned} \zeta(s) &:= \sum_{n=1}^{\infty} n^{-s}, \\ Z_n(s) &:= \sum'_{k \in \mathbb{Z}^n} |k|^{-s}, \\ \zeta_{p_1, \dots, p_n}(s) &:= \sum'_{k \in \mathbb{Z}^n} \frac{k_1^{p_1} \cdots k_n^{p_n}}{|k|^s}, \text{ for } p_i \in \mathbb{N}, \end{aligned}$$

where $\zeta(s)$ is the Riemann zeta function (see [56] or [36]).

By the symmetry $k \rightarrow -k$, it is clear that these functions ζ_{p_1, \dots, p_n} all vanish for odd values of p_i .

Let us now compute $\zeta_{0, \dots, 0, 1_i, 0, \dots, 0, 1_j, 0, \dots, 0}(s)$ in terms of $Z_n(s)$:

Since $\zeta_{0, \dots, 0, 1_i, 0, \dots, 0, 1_j, 0, \dots, 0}(s) = A_i(s) \delta_{ij}$, exchanging the components k_i and k_j , we get

$$\zeta_{0, \dots, 0, 1_i, 0, \dots, 0, 1_j, 0, \dots, 0}(s) = \frac{\delta_{ij}}{n} Z_n(s - 2).$$

Similarly,

$$\sum'_{\mathbb{Z}^n} \frac{k_1^2 k_2^2}{|k|^{s+8}} = \frac{1}{n(n-1)} Z_n(s+4) - \frac{1}{n-1} \sum'_{\mathbb{Z}^n} \frac{k_1^4}{|k|^{s+8}}$$

but it is difficult to write explicitly $\zeta_{p_1, \dots, p_n}(s)$ in terms of $Z_n(s-4)$ and other $Z_n(s-m)$ when at least four indices p_i are non zero.

When all p_i are even, $\zeta_{p_1, \dots, p_n}(s)$ is a nonzero series of fractions $\frac{P(k)}{|k|^s}$ where P is a homogeneous polynomial of degree $p_1 + \cdots + p_n$. Theorem 7.1 now gives us the following

Proposition 7.16. ζ_{p_1, \dots, p_n} has a meromorphic extension to the whole plane with a unique pole at $n + p_1 + \cdots + p_n$. This pole is simple and the residue at this pole is

$$\operatorname{Res}_{s=n+p_1+\dots+p_n} \zeta_{p_1, \dots, p_n}(s) = 2 \frac{\Gamma(\frac{p_1+1}{2}) \cdots \Gamma(\frac{p_n+1}{2})}{\Gamma(\frac{n+p_1+\dots+p_n}{2})} \quad (88)$$

when all p_i are even or this residue is zero otherwise.

In particular, for $n = 2$,

$$\operatorname{Res}_{s=0} \sum'_{k \in \mathbb{Z}^2} \frac{k_i k_j}{|k|^{s+4}} = \delta_{ij} \pi, \quad (89)$$

and for $n = 4$,

$$\begin{aligned} \operatorname{Res}_{s=0} \sum'_{k \in \mathbb{Z}^4} \frac{k_i k_j}{|k|^{s+6}} &= \delta_{ij} \frac{\pi^2}{2}, \\ \operatorname{Res}_{s=0} \sum'_{k \in \mathbb{Z}^4} \frac{k_i k_j k_l k_m}{|k|^{s+8}} &= (\delta_{ij} \delta_{lm} + \delta_{il} \delta_{jm} + \delta_{im} \delta_{jl}) \frac{\pi^2}{12}. \end{aligned} \quad (90)$$

Proof. Equation (88) follows from Theorem (7.1)

$$\operatorname{Res}_{s=n+p_1+\dots+p_n} \zeta_{p_1, \dots, p_n}(s) = \int_{k \in S^{n-1}} k_1^{p_1} \cdots k_n^{p_n} dS(k)$$

and standard formulae (see for instance [100, VIII,1;22]). Equation (89) is a straightforward consequence of Equation (88). Equation (90) can be checked for the cases $i = j \neq l = m$ and $i = j = l = m$. \square

Remark that $Z_n(s)$ is an Epstein zeta-function which is associated to the quadratic form $q(x) := x_1^2 + \dots + x_n^2$, so Z_n satisfies the following functional equation

$$Z_n(s) = \pi^{s-n/2} \Gamma(n/2 - s/2) \Gamma(s/2)^{-1} Z_n(n - s).$$

Since $\pi^{s-n/2} \Gamma(n/2 - s/2) \Gamma(s/2)^{-1} = 0$ for any negative even integer n and $Z_n(s)$ is meromorphic on \mathbb{C} with only one pole at $s = n$ with residue $2\pi^{n/2} \Gamma(n/2)^{-1}$ according to previous proposition, so we get $Z_n(0) = -1$. We have proved that

$$\operatorname{Res}_{s=0} Z_n(s + n) = 2\pi^{n/2} \Gamma(n/2)^{-1}, \quad (91)$$

$$Z_n(0) = -1. \quad (92)$$

There are many applications of Proposition 7.16 for instance in ζ -regularization, multiplicative anomalies or Casimir effect, see for instance [36].

7.4 Meromorphic continuation of a class of zeta functions

Let $n, q \in \mathbb{N}$, $q \geq 2$, and $p = (p_1, \dots, p_{q-1}) \in \mathbb{N}_0^{q-1}$.

Set $I := \{i \mid p_i \neq 0\}$ and assume that $I \neq \emptyset$ and

$$\mathcal{I} := \{\alpha = (\alpha_i)_{i \in I} \mid \forall i \in I \alpha_i = (\alpha_{i,1}, \dots, \alpha_{i,p_i}) \in \mathbb{N}_0^{p_i}\} = \prod_{i \in I} \mathbb{N}_0^{p_i}.$$

We will use in the sequel also the following notations:

- for $x = (x_1, \dots, x_t) \in \mathbb{R}^t$ recall that $|x|_1 = |x_1| + \dots + |x_t|$ and $|x| = \sqrt{x_1^2 + \dots + x_t^2}$;
- for all $\alpha = (\alpha_i)_{i \in I} \in \mathcal{I} = \prod_{i \in I} \mathbb{N}_0^{p_i}$,

$$|\alpha|_1 = \sum_{i \in I} |\alpha_i|_1 = \sum_{i \in I} \sum_{j=1}^{p_i} |\alpha_{i,j}| \quad \text{and} \quad \binom{1/2}{\alpha} = \prod_{i \in I} \binom{1/2}{\alpha_i} = \prod_{i \in I} \prod_{j=1}^{p_i} \binom{1/2}{\alpha_{i,j}}.$$

7.4.1 A family of polynomials

In this paragraph we define a family of polynomials which plays an important role later.

Consider first the variables:

- for X_1, \dots, X_n we set $X = (X_1, \dots, X_n)$;
 - for any $i = 1, \dots, 2q$, we consider the variables $Y_{i,1}, \dots, Y_{i,n}$ and set $Y_i := (Y_{i,1}, \dots, Y_{i,n})$ and $Y := (Y_1, \dots, Y_{2q})$;
 - for $Y = (Y_1, \dots, Y_{2q})$, we set for any $1 \leq j \leq q$, $\tilde{Y}_j := Y_1 + \dots + Y_j + Y_{q+1} + \dots + Y_{q+j}$.
- We define for all $\alpha = (\alpha_i)_{i \in I} \in \mathcal{I} = \prod_{i \in I} \mathbb{N}_0^{p_i}$ the polynomial

$$P_\alpha(X, Y) := \prod_{i \in I} \prod_{j=1}^{p_i} (2\langle X, \tilde{Y}_i \rangle + |\tilde{Y}_i|^2)^{\alpha_{i,j}}. \quad (93)$$

It is clear that $P_\alpha(X, Y) \in \mathbb{Z}[X, Y]$, $\deg_X P_\alpha \leq |\alpha|_1$ and $\deg_Y P_\alpha \leq 2|\alpha|_1$.

Let us fix a polynomial $Q \in \mathbb{R}[X_1, \dots, X_n]$ and note $d := \deg Q$. For $\alpha \in \mathcal{I}$, we want to expand $P_\alpha(X, Y) Q(X)$ in homogeneous polynomials in X and Y so defining

$$L(\alpha) := \{ \beta \in \mathbb{N}_0^{(2q+1)n} \mid |\beta|_1 - d_\beta \leq 2|\alpha|_1 \text{ and } d_\beta \leq |\alpha|_1 + d \}$$

where $d_\beta := \sum_1^n \beta_i$, we set

$$\binom{1/2}{\alpha} P_\alpha(X, Y) Q(X) =: \sum_{\beta \in L(\alpha)} c_{\alpha, \beta} X^\beta Y^\beta$$

where $c_{\alpha, \beta} \in \mathbb{R}$, $X^\beta := X_1^{\beta_1} \dots X_n^{\beta_n}$ and $Y^\beta := Y_{1,1}^{\beta_{n+1}} \dots Y_{2q,n}^{\beta_{(q+1)n}}$. By definition, X^β is a homogeneous polynomial of degree in X equals to d_β . We note

$$M_{\alpha, \beta}(Y) := c_{\alpha, \beta} Y^\beta.$$

7.4.2 Residues of a class of zeta functions

In this section we will prove the following result, used in Proposition 8.5 for the computation of the spectrum dimension of the noncommutative torus:

Theorem 7.17. (i) Let $\frac{1}{2\pi}\Theta$ be a badly approximable matrix, and $\tilde{a} \in \mathcal{S}((\mathbb{Z}^n)^{2q})$. Then

$$s \mapsto f(s) := \sum_{l \in [(\mathbb{Z}^n)^q]^2} \tilde{a}_l \sum'_{k \in \mathbb{Z}^n} \prod_{i=1}^{q-1} |k + \tilde{l}_i|^{p_i} |k|^{-s} Q(k) e^{ik \cdot \Theta \sum_1^q l_j}$$

has a meromorphic continuation to the whole complex plane \mathbb{C} with at most simple possible poles at the points $s = n + d + |p|_1 - m$ where $m \in \mathbb{N}_0$.

(ii) Let $m \in \mathbb{N}_0$ and set $I(m) := \{ (\alpha, \beta) \in \mathcal{I} \times \mathbb{N}_0^{(2q+1)n} \mid \beta \in L(\alpha) \text{ where we have taken } m = 2|\alpha|_1 - d_\beta + d \}$. Then $I(m)$ is a finite set and $s = n + d + |p|_1 - m$ is a pole of f if and only if

$$C(f, m) := \sum_{l \in Z} \tilde{a}_l \sum_{(\alpha, \beta) \in I(m)} M_{\alpha, \beta}(l) \int_{u \in S^{n-1}} u^\beta dS(u) \neq 0,$$

with $Z := \{ l \mid \sum_1^q l_j = 0 \}$ and the convention $\sum_\emptyset = 0$. In that case $s = n + d + |p|_1 - m$ is a simple pole of residue $\text{Res}_{s=n+d+|p|_1-m} f(s) = C(f, m)$.

In order to prove the theorem above we need the following

Lemma 7.18. *For all $N \in \mathbb{N}$ we have*

$$\prod_{i=1}^{q-1} |k + \tilde{l}_i|^{p_i} = \sum_{\alpha = (\alpha_i)_{i \in I} \in \prod_{i \in I} \{0, \dots, N\}^{p_i}} \binom{1/2}{\alpha} \frac{P_\alpha(k, l)}{|k|^{2|\alpha|_1 - |p|_1}} + \mathcal{O}_N(|k|^{p|1 - (N+1)/2})$$

uniformly in $k \in \mathbb{Z}^n$ and $l \in (\mathbb{Z}^n)^{2q}$ such that $|k| > U(l) := 36 (\sum_{i=1, i \neq q}^{2q-1} |l_i|)^4$.

Proof. For $i = 1, \dots, q-1$, we have uniformly in $k \in \mathbb{Z}^n$ and $l \in (\mathbb{Z}^n)^{2q}$ with $|k| > U(l)$,

$$\frac{|2\langle k, \tilde{l}_i \rangle + |\tilde{l}_i|^2|}{|k|^2} \leq \frac{\sqrt{U(l)}}{2|k|} < \frac{1}{2\sqrt{|k|}}. \quad (94)$$

In that case,

$$|k + \tilde{l}_i| = \left(|k|^2 + 2\langle k, \tilde{l}_i \rangle + |\tilde{l}_i|^2\right)^{1/2} = |k| \left(1 + \frac{2\langle k, \tilde{l}_i \rangle + |\tilde{l}_i|^2}{|k|^2}\right)^{1/2} = \sum_{u=0}^{\infty} \binom{1/2}{u} \frac{1}{|k|^{2u-1}} P_u^i(k, l)$$

where for all $i = 1, \dots, q-1$ and for all $u \in \mathbb{N}_0$,

$$P_u^i(k, l) := \left(2\langle k, \tilde{l}_i \rangle + |\tilde{l}_i|^2\right)^u,$$

with the convention $P_0^i(k, l) := 1$.

In particular $P_u^i(k, l) \in \mathbb{Z}[k, l]$, $\deg_k P_u^i \leq u$ and $\deg_l P_u^i \leq 2u$. Inequality (94) implies that for all $i = 1, \dots, q-1$ and for all $u \in \mathbb{N}$,

$$\frac{1}{|k|^{2u}} |P_u^i(k, l)| \leq \left(2\sqrt{|k|}\right)^{-u}$$

uniformly in $k \in \mathbb{Z}^n$ and $l \in (\mathbb{Z}^n)^{2q}$ such that $|k| > U(l)$.

Let $N \in \mathbb{N}$. We deduce from the previous that for any $k \in \mathbb{Z}^n$ and $l \in (\mathbb{Z}^n)^{2q}$ with $|k| > U(l)$ and for all $i = 1, \dots, q-1$, we have

$$\begin{aligned} |k + \tilde{l}_i| &= \sum_{u=0}^N \binom{1/2}{u} \frac{1}{|k|^{2u-1}} P_u^i(k, l) + \mathcal{O}\left(\sum_{u>N} |k| \binom{1/2}{u} (2\sqrt{|k|})^{-u}\right) \\ &= \sum_{u=0}^N \binom{1/2}{u} \frac{1}{|k|^{2u-1}} P_u^i(k, l) + \mathcal{O}_N\left(\frac{1}{|k|^{(N-1)/2}}\right). \end{aligned}$$

It follows that for any $N \in \mathbb{N}$, we have uniformly in $k \in \mathbb{Z}^n$ and $l \in (\mathbb{Z}^n)^{2q}$ with $|k| > U(l)$ and for all $i \in I$,

$$|k + \tilde{l}_i|^{p_i} = \sum_{\alpha_i \in \{0, \dots, N\}^{p_i}} \binom{1/2}{\alpha_i} \frac{1}{|k|^{2|\alpha_i|_1 - p_i}} P_{\alpha_i}^i(k, l) + \mathcal{O}_N\left(\frac{1}{|k|^{(N+1)/2 - p_i}}\right)$$

where $P_{\alpha_i}^i(k, l) = \prod_{j=1}^{p_i} P_{\alpha_{i,j}}^i(k, l)$ for all $\alpha_i = (\alpha_{i,1}, \dots, \alpha_{i,p_i}) \in \{0, \dots, N\}^{p_i}$ and

$$\prod_{i \in I} |k + \tilde{l}_i|^{p_i} = \sum_{\alpha = (\alpha_i) \in \prod_{i \in I} \{0, \dots, N\}^{p_i}} \binom{1/2}{\alpha} \frac{1}{|k|^{2|\alpha|_1 - |p|_1}} P_\alpha(k, l) + \mathcal{O}_N\left(\frac{1}{|k|^{(N+1)/2 - |p|_1}}\right)$$

where $P_\alpha(k, l) = \prod_{i \in I} P_{\alpha_i}^i(k, l) = \prod_{i \in I} \prod_{j=1}^{p_i} P_{\alpha_{i,j}}^i(k, l)$. \square

Proof of Theorem 7.17. (i) All $n, q, p = (p_1, \dots, p_{q-1})$ and $\tilde{a} \in \mathcal{S}((\mathbb{Z}^n)^{2q})$ are fixed as above and we define formally for any $l \in (\mathbb{Z}^n)^{2q}$

$$F(l, s) := \sum'_{k \in \mathbb{Z}^n} \prod_{i=1}^{q-1} |k + \tilde{l}_i|^{p_i} Q(k) e^{ik \cdot \Theta \sum_1^q l_j} |k|^{-s}. \quad (95)$$

Thus, still formally,

$$f(s) := \sum_{l \in (\mathbb{Z}^n)^{2q}} \tilde{a}_l F(l, s). \quad (96)$$

It is clear that $F(l, s)$ converges absolutely in the half plane $\{\sigma = \Re(s) > n + d + |p|_1\}$ where $d = \deg Q$.

Let $N \in \mathbb{N}$. Lemma 7.18 implies that for any $l \in (\mathbb{Z}^n)^{2q}$ and for $s \in \mathbb{C}$ such that $\sigma > n + |p|_1 + d$,

$$\begin{aligned} F(l, s) &= \sum'_{|k| \leq U(l)} \prod_{i=1}^{q-1} |k + \tilde{l}_i|^{p_i} Q(k) e^{ik \cdot \Theta \sum_1^q l_j} |k|^{-s} \\ &\quad + \sum_{\alpha = (\alpha_i)_{i \in I} \in \prod_{i \in I} \{0, \dots, N\}^{p_i}} \binom{1/2}{\alpha} \sum_{|k| > U(l)} \frac{1}{|k|^{s+2|\alpha|_1 - |p|_1}} P_\alpha(k, l) Q(k) e^{ik \cdot \Theta \sum_1^q l_j} + G_N(l, s). \end{aligned}$$

where $s \mapsto G_N(l, s)$ is a holomorphic function in the half-plane $D_N := \{\sigma > n + d + |p|_1 - \frac{N+1}{2}\}$ and verifies in it the bound $G_N(l, s) \ll_{N, \sigma} 1$ uniformly in l .

It follows that

$$F(l, s) = \sum_{\alpha = (\alpha_i)_{i \in I} \in \prod_{i \in I} \{0, \dots, N\}^{p_i}} H_\alpha(l, s) + R_N(l, s), \quad (97)$$

where

$$\begin{aligned} H_\alpha(l, s) &:= \sum'_{k \in \mathbb{Z}^n} \binom{1/2}{\alpha} \frac{1}{|k|^{s+2|\alpha|_1 - |p|_1}} P_\alpha(k, l) Q(k) e^{ik \cdot \Theta \sum_1^q l_j}, \\ R_N(l, s) &:= \sum'_{|k| \leq U(l)} \prod_{i=1}^{q-1} |k + \tilde{l}_i|^{p_i} Q(k) e^{ik \cdot \Theta \sum_1^q l_j} |k|^{-s} \\ &\quad - \sum'_{|k| \leq U(l)} \sum_{\alpha = (\alpha_i)_{i \in I} \in \prod_{i \in I} \{0, \dots, N\}^{p_i}} \binom{1/2}{\alpha} \frac{P_\alpha(k, l)}{|k|^{s+2|\alpha|_1 - |p|_1}} Q(k) e^{ik \cdot \Theta \sum_1^q l_j} + G_N(l, s). \end{aligned}$$

In particular there exists $A(N) > 0$ such that $s \mapsto R_N(l, s)$ extends holomorphically to the half-plane D_N and verifies in it the bound $R_N(l, s) \ll_{N, \sigma} 1 + |l|^{A(N)}$ uniformly in l .

Let us note formally

$$h_\alpha(s) := \sum_l \tilde{a}_l H_\alpha(l, s).$$

Equation (97) and $R_N(l, s) \ll_{N, \sigma} 1 + |l|^{A(N)}$ imply that

$$f(s) \sim_N \sum_{\alpha = (\alpha_i)_{i \in I} \in \prod_{i \in I} \{0, \dots, N\}^{p_i}} h_\alpha(s), \quad (98)$$

where \sim_N means modulo a holomorphic function in D_N .

Recall the decomposition $\binom{1/2}{\alpha} P_\alpha(k, l) Q(k) = \sum_{\beta \in L(\alpha)} M_{\alpha, \beta}(l) k^\beta$ and we decompose similarly $h_\alpha(s) = \sum_{\beta \in L(\alpha)} h_{\alpha, \beta}(s)$.

Theorem 7.5 now implies that for all $\alpha = (\alpha_i)_{i \in I} \in \prod_{i \in I} \{0, \dots, N\}^{p_i}$ and $\beta \in L(\alpha)$,

- the map $s \mapsto h_{\alpha, \beta}(s)$ has a meromorphic continuation to the whole complex plane \mathbb{C} with only one simple possible pole at $s = n + |p|_1 - 2|\alpha|_1 + d_\beta$,
- the residue at this point is equal to

$$\operatorname{Res}_{s=n+|p|_1-2|\alpha|_1+d_\beta} h_{\alpha, \beta}(s) = \sum_{l \in \mathcal{Z}} \tilde{a}_l M_{\alpha, \beta}(l) \int_{u \in S^{n-1}} u^\beta dS(u) \quad (99)$$

where $\mathcal{Z} := \{l \in (\mathbb{Z}^n)^{2q} : \sum_1^q l_j = 0\}$. If the right hand side is zero, $h_{\alpha, \beta}(s)$ is holomorphic on \mathbb{C} .

By (98), we deduce therefore that $f(s)$ has a meromorphic continuation on the halfplane D_N , with only simple possible poles in the set $\{n + |p|_1 + k : -2N|p|_1 \leq k \leq d\}$. Taking now $N \rightarrow \infty$ yields the result.

(ii) Let $m \in \mathbb{N}_0$ and set $I(m) := \{(\alpha, \beta) \in \mathcal{I} \times \mathbb{N}_0^{(2q+1)n} \mid \beta \in L(\alpha) \text{ and } m = 2|\alpha|_1 - d_\beta + d\}$. If $(\alpha, \beta) \in I(m)$, then $|\alpha|_1 \leq m$ and $|\beta|_1 \leq 3m + d$, so $I(m)$ is finite.

With a chosen N such that $2N|p|_1 + d > m$, we get by (98) and (99)

$$\operatorname{Res}_{s=n+d+|p|_1-m} f(s) = \sum_{l \in \mathcal{Z}} \tilde{a}_l \sum_{(\alpha, \beta) \in I(m)} M_{\alpha, \beta}(l) \int_{u \in S^{n-1}} u^\beta dS(u) = C(f, m)$$

with the convention $\sum_\emptyset = 0$. Thus, $n + d + |p|_1 - m$ is a pole of f if and only if $C(f, m) \neq 0$. \square

8 The noncommutative torus

The aim of this section is to compute the spectral action of the noncommutative torus. After the basic definitions, the result is presented in Theorem 8.13. Due to a fundamental appearance of small divisors, the number theory is involved via a Diophantine condition. As a consequence, the result which essentially says that the spectral action of the noncommutative torus coincide with the action of the ordinary torus (up few constants) is awfully technical and use the machinery of Section 7. A bunch of proofs are not given, but the essential lemmas are here: they show to the reader how life can be hard in noncommutative geometry!

Reference: [37].

8.1 Definition of the nc-torus

Let $C^\infty(\mathbb{T}_\Theta^n)$ be the smooth noncommutative n -torus associated to a non-zero skew-symmetric deformation matrix $\Theta \in M_n(\mathbb{R})$. It was introduced by Rieffel [95] and Connes [22] to generalize the n -torus \mathbb{T}^n .

This means that $C^\infty(\mathbb{T}_\Theta^n)$ is the algebra generated by n unitaries u_i , $i = 1, \dots, n$ subject to the relations

$$u_l u_j = e^{i\Theta_{lj}} u_j u_l, \quad (100)$$

and with Schwartz coefficients: an element $a \in C^\infty(\mathbb{T}_\Theta^n)$ can be written as $a = \sum_{k \in \mathbb{Z}^n} a_k U_k$, where $\{a_k\} \in \mathcal{S}(\mathbb{Z}^n)$ with the Weyl elements defined by

$$U_k := e^{-\frac{i}{2}k \cdot \chi^k} u_1^{k_1} \dots u_n^{k_n},$$

$k \in \mathbb{Z}^n$, relation (100) reads

$$U_k U_q = e^{-\frac{i}{2}k \cdot \Theta q} U_{k+q}, \text{ and } U_k U_q = e^{-ik \cdot \Theta q} U_q U_k \quad (101)$$

where χ is the matrix restriction of Θ to its upper triangular part. Thus unitary operators U_k satisfy

$$U_k^* = U_{-k} \text{ and } [U_k, U_l] = -2i \sin(\frac{1}{2}k \cdot \Theta l) U_{k+l}.$$

Let τ be the trace on $C^\infty(\mathbb{T}_\Theta^n)$ defined by

$$\tau\left(\sum_{k \in \mathbb{Z}^n} a_k U_k\right) := a_0$$

and \mathcal{H}_τ be the GNS Hilbert space obtained by completion of $C^\infty(\mathbb{T}_\Theta^n)$ with respect of the norm induced by the scalar product

$$\langle a, b \rangle := \tau(a^* b).$$

On $\mathcal{H}_\tau = \{ \sum_{k \in \mathbb{Z}^n} a_k U_k \mid \{a_k\}_k \in l^2(\mathbb{Z}^n) \}$, we consider the left and right regular representations of $C^\infty(\mathbb{T}_\Theta^n)$ by bounded operators, that we denote respectively by $L(\cdot)$ and $R(\cdot)$.

Let also δ_μ , $\mu \in \{1, \dots, n\}$, be the n (pairwise commuting) canonical derivations, defined by

$$\delta_\mu(U_k) := ik_\mu U_k. \quad (102)$$

We need to fix notations: let

$$\mathcal{A}_\Theta := C^\infty(\mathbb{T}_\Theta^n) \text{ acting on } \mathcal{H} := \mathcal{H}_\tau \otimes \mathbb{C}^{2^m}$$

with $n = 2m$ or $n = 2m + 1$ (i.e., $m = \lfloor \frac{n}{2} \rfloor$ is the integer part of $\frac{n}{2}$), the square integrable sections of the trivial spin bundle over \mathbb{T}^m .

Each element of \mathcal{A}_Θ is represented on \mathcal{H} as $L(a) \otimes 1_{2^m}$. The Tomita conjugation

$$J_0(a) := a^*$$

satisfies $[J_0, \delta_\mu] = 0$ and we define

$$J := J_0 \otimes C_0$$

where C_0 is an operator on \mathbb{C}^{2^m} . The Dirac-like operator is given by

$$\mathcal{D} := -i \delta_\mu \otimes \gamma^\mu, \tag{103}$$

where we use hermitian Dirac matrices γ . It is defined and symmetric on the dense subset of \mathcal{H} given by $C^\infty(\mathbb{T}_\Theta^n) \otimes \mathbb{C}^{2^m}$. We still note \mathcal{D} its selfadjoint extension. This implies

$$C_0 \gamma^\alpha = -\varepsilon \gamma^\alpha C_0, \tag{104}$$

and

$$\mathcal{D} U_k \otimes e_i = k_\mu U_k \otimes \gamma^\mu e_i,$$

where (e_i) is the canonical basis of \mathbb{C}^{2^m} . Moreover, $C_0^2 = \pm 1_{2^m}$ depending on the parity of m . Finally, one introduces the chirality, which in the even case is

$$\chi := id \otimes (-i)^m \gamma^1 \cdots \gamma^n.$$

This yields a spectral triple:

Theorem 8.1. *The 5-tuple $(\mathcal{A}_\Theta, \mathcal{H}, \mathcal{D}, J, \chi)$ is a real regular spectral triple of dimension n . It satisfies the finiteness and orientability conditions of Definition 5.2. It is n -summable and its KO -dimension is also n .*

We do not give a proof since most of its arguments will be emphasized in this section; see however [25, 50] for a specific proof.

For instance, we prove in Proposition 8.5 that this triple has simple dimension spectrum when Θ is badly approximable (see Definition 7.4).

The perturbed Dirac operator $V_u \mathcal{D} V_u^*$ by the unitary

$$V_u := \left(L(u) \otimes 1_{2^m} \right) J \left(L(u) \otimes 1_{2^m} \right) J^{-1},$$

defined for every unitary $u \in \mathcal{A}$, $uu^* = u^*u = U_0$, must satisfy condition $J\mathcal{D} = \epsilon \mathcal{D}J$ (which is equivalent to \mathcal{H} being endowed with a structure of \mathcal{A}_Θ -bimodule). This yields the necessity of a symmetrized covariant Dirac operator

$$\mathcal{D}_A := \mathcal{D} + A + \epsilon J A J^{-1}$$

since $V_u \mathcal{D} V_u^* = \mathcal{D}_{L(u) \otimes 1_{2^m} [\mathcal{D}, L(u^*) \otimes 1_{2^m}]}$: in fact, for $a \in \mathcal{A}_\Theta$, using $J_0 L(a) J_0^{-1} = R(a^*)$, we get

$$\epsilon J \left(L(a) \otimes \gamma^\alpha \right) J^{-1} = -R(a^*) \otimes \gamma^\alpha$$

and that the representation L and the anti-representation R are \mathbb{C} -linear, commute and satisfy

$$[\delta_\alpha, L(a)] = L(\delta_\alpha a), \quad [\delta_\alpha, R(a)] = R(\delta_\alpha a).$$

This induces some covariance property for the Dirac operator: one checks that for all $k \in \mathbb{Z}^n$,

$$L(U_k) \otimes 1_{2^m} [\mathcal{D}, L(U_k^*) \otimes 1_{2^m}] = 1 \otimes (-k_\mu \gamma^\mu), \quad (105)$$

so with (104), we get $U_k[\mathcal{D}, U_k^*] + \epsilon J U_k[\mathcal{D}, U_k^*] J^{-1} = 0$ and

$$V_{U_k} \mathcal{D} V_{U_k}^* = \mathcal{D} = \mathcal{D}_{L(U_k) \otimes 1_{2^m} [\mathcal{D}, L(U_k^*) \otimes 1_{2^m}]}. \quad (106)$$

Moreover, we get the gauge transformation (see Lemma 5.13):

$$V_u \mathcal{D}_A V_u^* = \mathcal{D}_{\gamma_u(A)} \quad (107)$$

where the gauged transform one-form of A is

$$\gamma_u(A) := u[\mathcal{D}, u^*] + u A u^*, \quad (108)$$

with the shorthand $L(u) \otimes 1_{2^m} \rightarrow u$. As a consequence, the spectral action is gauge invariant:

$$\mathcal{S}(\mathcal{D}_A, f, \Lambda) = \mathcal{S}(\mathcal{D}_{\gamma_u(A)}, f, \Lambda).$$

An arbitrary selfadjoint one-form $A \in \Omega_D^1(\mathcal{A})$, can be written as

$$A = L(-i A_\alpha) \otimes \gamma^\alpha, \quad A_\alpha = -A_\alpha^* \in \mathcal{A}_\Theta, \quad (109)$$

thus

$$\mathcal{D}_A = -i \left(\delta_\alpha + L(A_\alpha) - R(A_\alpha) \right) \otimes \gamma^\alpha. \quad (110)$$

Defining

$$\tilde{A}_\alpha := L(A_\alpha) - R(A_\alpha),$$

we get $\mathcal{D}_A^2 = -g^{\alpha_1 \alpha_2} (\delta_{\alpha_1} + \tilde{A}_{\alpha_1}) (\delta_{\alpha_2} + \tilde{A}_{\alpha_2}) \otimes 1_{2^m} - \frac{1}{2} \Omega_{\alpha_1 \alpha_2} \otimes \gamma^{\alpha_1 \alpha_2}$ where

$$\begin{aligned} \gamma^{\alpha_1 \alpha_2} &:= \frac{1}{2} (\gamma^{\alpha_1} \gamma^{\alpha_2} - \gamma^{\alpha_2} \gamma^{\alpha_1}), \\ \Omega_{\alpha_1 \alpha_2} &:= [\delta_{\alpha_1} + \tilde{A}_{\alpha_1}, \delta_{\alpha_2} + \tilde{A}_{\alpha_2}] = L(F_{\alpha_1 \alpha_2}) - R(F_{\alpha_1 \alpha_2}) \end{aligned}$$

with

$$F_{\alpha_1 \alpha_2} := \delta_{\alpha_1}(A_{\alpha_2}) - \delta_{\alpha_2}(A_{\alpha_1}) + [A_{\alpha_1}, A_{\alpha_2}]. \quad (111)$$

In summary,

$$\begin{aligned} \mathcal{D}_A^2 &= -\delta^{\alpha_1 \alpha_2} \left(\delta_{\alpha_1} + L(A_{\alpha_1}) - R(A_{\alpha_1}) \right) \left(\delta_{\alpha_2} + L(A_{\alpha_2}) - R(A_{\alpha_2}) \right) \otimes 1_{2^m} \\ &\quad - \frac{1}{2} \left(L(F_{\alpha_1 \alpha_2}) - R(F_{\alpha_1 \alpha_2}) \right) \otimes \gamma^{\alpha_1 \alpha_2}. \end{aligned} \quad (112)$$

8.2 Kernels and dimension spectrum

We now compute the kernel of the perturbed Dirac operator:

Proposition 8.2. (i) $\text{Ker } \mathcal{D} = U_0 \otimes \mathbb{C}^{2^m}$, so $\dim \text{Ker } \mathcal{D} = 2^m$.

(ii) For any selfadjoint one-form A , $\text{Ker } \mathcal{D} \subseteq \text{Ker } \mathcal{D}_A$.

(iii) For any unitary $u \in \mathcal{A}$, $\text{Ker } \mathcal{D}_{\gamma_u(A)} = V_u \text{Ker } \mathcal{D}_A$.

Proof. (i) Let $\psi = \sum_{k,j} c_{k,j} U_k \otimes e_j \in \text{Ker } \mathcal{D}$. Thus, $0 = \mathcal{D}^2 \psi = \sum_{k,i} c_{k,j} |k|^2 U_k \otimes e_j$ which entails that $c_{k,j} |k|^2 = 0$ for any $k \in \mathbb{Z}^n$ and $1 \leq j \leq 2^m$. The result follows.

(ii) Let $\psi \in \text{Ker } \mathcal{D}$. So, $\psi = U_0 \otimes v$ with $v \in \mathbb{C}^{2^m}$ and from (110), we get

$$\mathcal{D}_A \psi = \mathcal{D} \psi + (A + \epsilon J A J^{-1}) \psi = (A + \epsilon J A J^{-1}) \psi = -i[A_\alpha, U_0] \otimes \gamma^\alpha v = 0$$

since U_0 is the unit of the algebra, which proves that $\psi \in \text{Ker } \mathcal{D}_A$.

(iii) This is a direct consequence of (107). \square

Corollary 8.3. Let A be a selfadjoint one-form. Then $\text{Ker } \mathcal{D}_A = \text{Ker } \mathcal{D}$ in the following cases:

(i) $A = A_u := L(u) \otimes 1_{2^m}[\mathcal{D}, L(u^*) \otimes 1_{2^m}]$ when u is a unitary in \mathcal{A} .

(ii) $\|A\| < \frac{1}{2}$.

(iii) The matrix $\frac{1}{2\pi} \Theta$ has only integral coefficients.

Proof. (i) This follows from previous result because $V_u(U_0 \otimes v) = U_0 \otimes v$ for any $v \in \mathbb{C}^{2^m}$.

(ii) Let $\psi = \sum_{k,j} c_{k,j} U_k \otimes e_j$ be in $\text{Ker } \mathcal{D}_A$ (so $\sum_{k,j} |c_{k,j}|^2 < \infty$) and $\phi := \sum_j c_{0,j} U_0 \otimes e_j$. Thus $\psi' := \psi - \phi \in \text{Ker } \mathcal{D}_A$ since $\phi \in \text{Ker } \mathcal{D} \subseteq \text{Ker } \mathcal{D}_A$ and

$$\| \sum_{0 \neq k \in \mathbb{Z}^n, j} c_{k,j} k_\alpha U_k \otimes \gamma^\alpha e_j \|^2 = \| \mathcal{D} \psi' \|^2 = \| -(A + \epsilon J A J^{-1}) \psi' \|^2 \leq 4 \|A\|^2 \| \psi' \|^2 < \| \psi' \|^2.$$

Defining $X_k := \sum_\alpha k_\alpha \gamma_\alpha$, $X_k^2 = \sum_\alpha |k_\alpha|^2 1_{2^m}$ is invertible and the vectors $\{U_k \otimes X_k e_j\}_{0 \neq k \in \mathbb{Z}^n, j}$ are orthogonal in \mathcal{H} , so

$$\sum_{0 \neq k \in \mathbb{Z}^n, j} \left(\sum_\alpha |k_\alpha|^2 \right) |c_{k,j}|^2 < \sum_{0 \neq k \in \mathbb{Z}^n, j} |c_{k,j}|^2$$

which is possible only if $c_{k,j} = 0$, $\forall k, j$ that is $\psi' = 0$ and $\psi = \phi \in \text{Ker } \mathcal{D}$.

(iii) This is a consequence of the fact that the algebra is commutative, thus the arguments of (64) apply and $\tilde{A} = 0$. \square

Note that if $\tilde{A}_u := A_u + \epsilon J A_u J^{-1}$, then by (105), $\tilde{A}_{U_k} = 0$ for all $k \in \mathbb{Z}^n$ and $\|A_{U_k}\| = |k|$, but for an arbitrary unitary $u \in \mathcal{A}$, $\tilde{A}_u \neq 0$ so $\mathcal{D}_{A_u} \neq \mathcal{D}$.

Naturally the above result is also a direct consequence of the fact that the eigenspace of an isolated eigenvalue of an operator is not modified by small perturbations. However, it is interesting to compute the last result directly to emphasize the difficulty of the general case:

Let $\psi = \sum_{l \in \mathbb{Z}^n, 1 \leq j \leq 2^m} c_{l,j} U_l \otimes e_j \in \text{Ker } \mathcal{D}_A$, so $\sum_{l \in \mathbb{Z}^n, 1 \leq j \leq 2^m} |c_{l,j}|^2 < \infty$. We have to show that $\psi \in \text{Ker } \mathcal{D}$ that is $c_{l,j} = 0$ when $l \neq 0$.

Taking the scalar product of $\langle U_k \otimes e_i |$ with

$$0 = \mathcal{D}_A \psi = \sum_{l, \alpha, j} c_{l,j} (l^\alpha U_l - i[A_\alpha, U_l]) \otimes \gamma^\alpha e_j,$$

we obtain

$$0 = \sum_{l, \alpha, j} c_{l, j} \left(l^\alpha \delta_{k, l} - i \langle U_k, [A_\alpha, U_l] \rangle \right) \langle e_i, \gamma^\alpha e_j \rangle.$$

If $A_\alpha = \sum_{\alpha, l} a_{\alpha, l} U_l \otimes \gamma^\alpha$ with $\{a_{\alpha, l}\}_l \in \mathcal{S}(\mathbb{Z}^n)$, note that $[U_l, U_m] = -2i \sin(\frac{1}{2}l \cdot \Theta m) U_{l+m}$ and

$$\langle U_k, [A_\alpha, U_l] \rangle = \sum_{l' \in \mathbb{Z}^n} a_{\alpha, l'} (-2i \sin(\frac{1}{2}l' \cdot \Theta l)) \langle U_k, U_{l'+l} \rangle = -2i a_{\alpha, k-l} \sin(\frac{1}{2}k \cdot \Theta l).$$

Thus

$$0 = \sum_{l \in \mathbb{Z}^n} \sum_{\alpha=1}^n \sum_{j=1}^{2^m} c_{l, j} \left(l^\alpha \delta_{k, l} - 2a_{\alpha, k-l} \sin(\frac{1}{2}k \cdot \Theta l) \right) \langle e_i, \gamma^\alpha e_j \rangle, \quad \forall k \in \mathbb{Z}^n, \forall i, 1 \leq i \leq 2^m. \quad (113)$$

We conjecture that $\text{Ker } \mathcal{D} = \text{Ker } \mathcal{D}_A$ at least for generic Θ 's: the constraints (113) should imply $c_{l, j} = 0$ for all j and all $l \neq 0$ meaning $\psi \in \text{Ker } \mathcal{D}$. When $\frac{1}{2\pi}\Theta$ has only integer coefficients, the sin part of these constraints disappears giving the result.

We will use freely the notation (49) about the difference between \mathcal{D} and D .

Lemma 8.4. *If $\frac{1}{2\pi}\Theta$ is badly approximable (see Definition 7.4), $Sp(C^\infty(\mathbb{T}_\Theta^n), \mathcal{H}, \mathcal{D}) = \mathbb{Z}$ and all these poles are simple.*

Proof. Let $B \in \mathcal{D}(\mathcal{A})$ and $p \in \mathbb{N}_0$. Suppose that B is of the form

$$B = a_r b_r \mathcal{D}^{q_{r-1}} |D|^{p_{r-1}} a_{r-1} b_{r-1} \cdots \mathcal{D}^{q_1} |D|^{p_1} a_1 b_1$$

where $r \in \mathbb{N}$, $a_i \in \mathcal{A}$, $b_i \in \mathcal{J}\mathcal{A}\mathcal{J}^{-1}$, $q_i, p_i \in \mathbb{N}_0$. We note $a_i =: \sum_l a_{i, l} U_l$ and $b_i =: \sum_l b_{i, l} U_l$. With the shorthand $k_{\mu_1, \mu_{q_i}} := k_{\mu_1} \cdots k_{\mu_{q_i}}$ and $\gamma^{\mu_1, \mu_{q_i}} = \gamma^{\mu_1} \cdots \gamma^{\mu_{q_i}}$, we get

$$\mathcal{D}^{q_1} |D|^{p_1} a_1 b_1 U_k \otimes e_j = \sum_{l_1, l'_1} a_{1, l_1} b_{1, l'_1} U_{l_1} U_k U_{l'_1} |k + l_1 + l'_1|^{p_1} (k + l_1 + l'_1)_{\mu_1, \mu_{q_1}} \otimes \gamma^{\mu_1, \mu_{q_1}} e_j$$

which gives, after r iterations,

$$B U_k \otimes e_j = \sum_{l, l'} \tilde{a}_l \tilde{b}_{l'} U_{l_r} \cdots U_{l_1} U_k U_{l'_1} \cdots U_{l'_r} \prod_{i=1}^{r-1} |k + \hat{l}_i + \hat{l}'_i|^{p_i} (k + \hat{l}_i + \hat{l}'_i)_{\mu_1^i, \mu_{q_i}^i} \otimes \gamma^{\mu_1^{r-1}, \mu_{q_{r-1}}^{r-1}} \cdots \gamma^{\mu_1^1, \mu_{q_1}^1} e_j$$

where $\tilde{a}_l := a_{1, l_1} \cdots a_{r, l_r}$ and $\tilde{b}_{l'} := b_{1, l'_1} \cdots b_{r, l'_r}$.

Let us note $F_\mu(k, l, l') := \prod_{i=1}^{r-1} |k + \hat{l}_i + \hat{l}'_i|^{p_i} (k + \hat{l}_i + \hat{l}'_i)_{\mu_1^i, \mu_{q_i}^i}$ and $\gamma^\mu := \gamma^{\mu_1^{r-1}, \mu_{q_{r-1}}^{r-1}} \cdots \gamma^{\mu_1^1, \mu_{q_1}^1}$. Thus, with the shortcut

\sim_c meaning modulo a constant function towards the variable s ,

$$\text{Tr} \left(B |D|^{-2p-s} \right) \sim_c \sum_k' \sum_{l, l'} \tilde{a}_l \tilde{b}_{l'} \tau \left(U_{-k} U_{l_r} \cdots U_{l_1} U_k U_{l'_1} \cdots U_{l'_r} \right) \frac{F_\mu(k, l, l')}{|k|^{s+2p}} \text{Tr}(\gamma^\mu).$$

Since $U_{l_r} \cdots U_{l_1} U_k = U_k U_{l_r} \cdots U_{l_1} e^{-i \sum_1^r l_i \cdot \Theta k}$ we get

$$\tau \left(U_{-k} U_{l_r} \cdots U_{l_1} U_k U_{l'_1} \cdots U_{l'_r} \right) = \delta_{\sum_1^r l_i + l'_i, 0} e^{i\phi(l, l')} e^{-i \sum_1^r l_i \cdot \Theta k}$$

where ϕ is a real valued function. Thus,

$$\begin{aligned} \mathrm{Tr} \left(B|D|^{-2p-s} \right) &\sim_c \sum_k' \sum_{l,l'} e^{i\phi(l,l')} \delta_{\sum_1^r l_i+l'_i,0} \tilde{a}_l \tilde{b}_{l'} \frac{F_\mu(k,l,l') e^{-i \sum_1^r l_i \cdot \Theta k}}{|k|^{s+2p}} \mathrm{Tr}(\gamma^\mu) \\ &\sim_c f_\mu(s) \mathrm{Tr}(\gamma^\mu). \end{aligned}$$

The function $f_\mu(s)$ can be decomposed as a linear combination of zeta function of type described in Theorem 7.17 (or, if $r = 1$ or all the p_i are zero, in Theorem 7.5). Thus, $s \mapsto \mathrm{Tr} \left(B|D|^{-2p-s} \right)$ has only poles in \mathbb{Z} and each pole is simple. Finally, by linearity, we get the result. \square

The dimension spectrum of the noncommutative torus is simple:

Proposition 8.5. (i) *If $\frac{1}{2\pi}\Theta$ is badly approximable, the spectrum dimension of the spectral triple $(C^\infty(\mathbb{T}_\Theta^n), \mathcal{H}, \mathcal{D})$ is equal to the set $\{n - k : k \in \mathbb{N}_0\}$ and all these poles are simple.*

(ii) $\zeta_D(0) = 0$.

Proof. (i) Lemma 8.4 and Remark 5.9.

(ii) $\zeta_D(s) = \sum_{k \in \mathbb{Z}^n} \sum_{1 \leq j \leq 2^m} \langle U_k \otimes e_j, |D|^{-s} U_k \otimes e_j \rangle = 2^m (\sum_{k \in \mathbb{Z}^n}' \frac{1}{|k|^s} + 1) = 2^m (Z_n(s) + 1)$. By (92), we get the result. \square

We have computed $\zeta_D(0)$ relatively easily but the main difficulty of the present Section is precisely to calculate $\zeta_{D_A}(0)$.

8.3 Noncommutative integral computations

We fix a self-adjoint one-form A on the noncommutative torus of dimension n .

Proposition 8.6. *If $\frac{1}{2\pi}\Theta$ is badly approximable, then the first elements of the spectral action expansion (74) are given by*

$$\begin{aligned} \int |D_A|^{-n} &= \int |D|^{-n} = 2^{m+1} \pi^{n/2} \Gamma(\frac{n}{2})^{-1}. \\ \int |D_A|^{-n+k} &= 0 \text{ for } k \text{ odd.} \\ \int |D_A|^{-n+2} &= 0. \end{aligned}$$

We need a few technical lemmas:

Lemma 8.7. *On the noncommutative torus, for any $t \in \mathbb{R}$,*

$$\int \tilde{A} \mathcal{D} |D|^{-t} = \int \mathcal{D} \tilde{A} |D|^{-t} = 0.$$

Proof. Using notations of (109), we have

$$\begin{aligned} \mathrm{Tr}(\tilde{A} \mathcal{D} |D|^{-s}) &\sim_c \sum_j \sum_k' \langle U_k \otimes e_j, -ik_\mu |k|^{-s} [A_\alpha, U_k] \otimes \gamma^\alpha \gamma^\mu e_j \rangle \\ &\sim_c -i \mathrm{Tr}(\gamma^\alpha \gamma^\mu) \sum_k' k_\mu |k|^{-s} \langle U_k, [A_\alpha, U_k] \rangle = 0 \end{aligned}$$

since $\langle U_k, [A_\alpha, U_k] \rangle = 0$. Similarly

$$\begin{aligned} \mathrm{Tr}(\mathcal{D} \tilde{A} |D|^{-s}) &\sim_c \sum_j \sum_k' \langle U_k \otimes e_j, |k|^{-s} \sum_l a_{\alpha,l} 2 \sin \frac{k \cdot \Theta l}{2} (l+k)_\mu U_{l+k} \otimes \gamma^\mu \gamma^\alpha e_j \rangle \\ &\sim_c 2 \mathrm{Tr}(\gamma^\mu \gamma^\alpha) \sum_k' \sum_l a_{\alpha,l} \sin \frac{k \cdot \Theta l}{2} (l+k)_\mu |k|^{-s} \langle U_k, U_{l+k} \rangle = 0. \end{aligned} \quad \square$$

Any element h in the algebra generated by \mathcal{A} and $[\mathcal{D}, \mathcal{A}]$ can be written as a linear combination of terms of the form $a_1^{p_1} \cdots a_n^{p_r}$ where a_i are elements of \mathcal{A} or $[\mathcal{D}, \mathcal{A}]$. Such a term can be written as a series $b := \sum a_{1,\alpha_1,l_1} \cdots a_{q,\alpha_q,l_q} U_{l_1} \cdots U_{l_q} \otimes \gamma^{\alpha_1} \cdots \gamma^{\alpha_q}$ where a_{i,α_i} are Schwartz sequences and when $a_i =: \sum_l a_l U_l \in \mathcal{A}$, we set $a_{i,\alpha,l} = a_{i,l}$ with $\gamma^\alpha = 1$. We define

$$L(b) := \tau \left(\sum_l a_{1,\alpha_1,l_1} \cdots a_{q,\alpha_q,l_q} U_{l_1} \cdots U_{l_q} \right) \text{Tr}(\gamma^{\alpha_1} \cdots \gamma^{\alpha_q}).$$

By linearity, L is defined as a linear form on the whole algebra generated by \mathcal{A} and $[\mathcal{D}, \mathcal{A}]$.

Lemma 8.8. *If h is an element of the algebra generated by \mathcal{A} and $[\mathcal{D}, \mathcal{A}]$,*

$$\text{Tr} \left(h |D|^{-s} \right) \sim_c L(h) Z_n(s).$$

In particular, $\text{Tr} \left(h |D|^{-s} \right)$ has at most one pole at $s = n$.

Proof. We get with b of the form $\sum a_{1,\alpha_1,l_1} \cdots a_{q,\alpha_q,l_q} U_{l_1} \cdots U_{l_q} \otimes \gamma^{\alpha_1} \cdots \gamma^{\alpha_q}$,

$$\begin{aligned} \text{Tr} \left(b |D|^{-s} \right) &\sim_c \sum'_{k \in \mathbb{Z}^n} \langle U_k, \sum_l a_{1,\alpha_1,l_1} \cdots a_{q,\alpha_q,l_q} U_{l_1} \cdots U_{l_q} U_k \rangle \text{Tr}(\gamma^{\alpha_1} \cdots \gamma^{\alpha_q}) |k|^{-s} \\ &\sim_c \tau \left(\sum_l a_{1,\alpha_1,l_1} \cdots a_{q,\alpha_q,l_q} U_{l_1} \cdots U_{l_q} \right) \text{Tr}(\gamma^{\alpha_1} \cdots \gamma^{\alpha_q}) Z_n(s) = L(b) Z_n(s). \end{aligned}$$

The results follows now from linearity of the trace. \square

Lemma 8.9. *If $\frac{1}{2\pi}\Theta$ is badly approximable, the function $s \mapsto \text{Tr} \left(\varepsilon JAJ^{-1}A |D|^{-s} \right)$ extends meromorphically on the whole plane with only one possible pole at $s = n$. Moreover, this pole is simple and*

$$\text{Res}_{s=n} \text{Tr} \left(\varepsilon JAJ^{-1}A |D|^{-s} \right) = a_{\alpha,0} a_0^\alpha 2^{m+1} \pi^{n/2} \Gamma(n/2)^{-1}.$$

Proof. With $A = L(-iA_\alpha) \otimes \gamma^\alpha$, we get $\varepsilon JAJ^{-1} = R(iA_\alpha) \otimes \gamma^\alpha$, and by multiplication $\varepsilon JAJ^{-1}A = R(A_\beta)L(A_\alpha) \otimes \gamma^\beta \gamma^\alpha$. Thus,

$$\begin{aligned} \text{Tr} \left(\varepsilon JAJ^{-1}A |D|^{-s} \right) &\sim_c \sum'_{k \in \mathbb{Z}^n} \langle U_k, A_\alpha U_k A_\beta \rangle |k|^{-s} \text{Tr}(\gamma^\beta \gamma^\alpha) \\ &\sim_c \sum'_{k \in \mathbb{Z}^n} \sum_l a_{\alpha,l} a_{\beta,-l} e^{ik \cdot \Theta l} |k|^{-s} \text{Tr}(\gamma^\beta \gamma^\alpha) \\ &\sim_c 2^m \sum'_{k \in \mathbb{Z}^n} \sum_l a_{\alpha,l} a_{-l}^\alpha e^{ik \cdot \Theta l} |k|^{-s}. \end{aligned}$$

Theorem 7.5 (ii) entails that $\sum'_{k \in \mathbb{Z}^n} \sum_l a_{\alpha,l} a_{-l}^\alpha e^{ik \cdot \Theta l} |k|^{-s}$ extends meromorphically on the whole plane \mathbb{C} with only one possible pole at $s = n$. Moreover, this pole is simple and we have

$$\text{Res}_{s=n} \sum'_{k \in \mathbb{Z}^n} \sum_l a_{\alpha,l} a_{-l}^\alpha e^{ik \cdot \Theta l} |k|^{-s} = a_{\alpha,0} a_0^\alpha \text{Res}_{s=n} Z_n(s).$$

Equation (91) now gives the result. \square

Lemma 8.10. *If $\frac{1}{2\pi}\Theta$ is badly approximable, then for any $t \in \mathbb{R}$,*

$$\int X |D|^{-t} = \delta_{t,n} 2^{m+1} \left(- \sum_l a_{\alpha,l} a_{-l}^\alpha + a_{\alpha,0} a_0^\alpha \right) 2\pi^{n/2} \Gamma(n/2)^{-1}.$$

where $X = \tilde{A}\mathcal{D} + \mathcal{D}\tilde{A} + \tilde{A}^2$ and $A =: -i \sum_l a_{\alpha,l} U_l \otimes \gamma^\alpha$.

Proof. By Lemma 8.7, we get $\int X|D|^{-t} = \text{Res}_{s=0} \text{Tr}(\tilde{A}^2|D|^{-s-t})$. Since A and εJAJ^{-1} commute, we have $\tilde{A}^2 = A^2 + JA^2J^{-1} + 2\varepsilon JAJ^{-1}A$. Thus,

$$\text{Tr}(\tilde{A}^2|D|^{-s-t}) = \text{Tr}(A^2|D|^{-s-t}) + \text{Tr}(JA^2J^{-1}|D|^{-s-t}) + 2\text{Tr}(\varepsilon JAJ^{-1}A|D|^{-s-t}).$$

Since $|D|$ and J commute, we have with Lemma 8.8,

$$\text{Tr}(\tilde{A}^2|D|^{-s-t}) \sim_c 2L(A^2)Z_n(s+t) + 2\text{Tr}(\varepsilon JAJ^{-1}A|D|^{-s-t}).$$

Thus Lemma 8.9 entails that $\text{Tr}(\tilde{A}^2|D|^{-s-t})$ is holomorphic at 0 if $t \neq n$. When $t = n$,

$$\text{Res}_{s=0} \text{Tr}(\tilde{A}^2|D|^{-s-t}) = 2^{m+1} \left(-\sum_l a_{\alpha,l} a_{-l}^\alpha + a_{\alpha,0} a_0^\alpha \right) 2\pi^{n/2} \Gamma(n/2)^{-1}, \quad (114)$$

which gives the result. \square

Lemma 8.11. *If $\frac{1}{2\pi}\Theta$ is badly approximable, then*

$$\int \tilde{A}\mathcal{D}\tilde{A}\mathcal{D}|D|^{-2-n} = -\frac{n-2}{n} \int \tilde{A}^2|D|^{-n}.$$

Proof. With $\mathcal{D}J = \varepsilon J\mathcal{D}$, we get

$$\int \tilde{A}\mathcal{D}\tilde{A}\mathcal{D}|D|^{-2-n} = 2\int A\mathcal{D}A\mathcal{D}|D|^{-2-n} + 2\int \varepsilon JAJ^{-1}\mathcal{D}A\mathcal{D}|D|^{-2-n}.$$

Let us first compute $\int A\mathcal{D}A\mathcal{D}|D|^{-2-n}$. We have, with $A =: -iL(A_\alpha) \otimes \gamma^\alpha =: -i \sum_l a_{\alpha,l} U_l \otimes \gamma^\alpha$,

$$\text{Tr}(A\mathcal{D}A\mathcal{D}|D|^{-s-2-n}) \sim_c -\sum_k' \sum_{l_1, l_2} a_{\alpha_2, l_2} a_{\alpha_1, l_1} \tau(U_{-k} U_{l_2} U_{l_1} U_k) \frac{k_{\mu_1} (k+l_1)_{\mu_2}}{|k|^{s+2+n}} \text{Tr}(\gamma^{\alpha, \mu})$$

where $\gamma^{\alpha, \mu} := \gamma^{\alpha_2} \gamma^{\mu_2} \gamma^{\alpha_1} \gamma^{\mu_1}$. Thus,

$$\int A\mathcal{D}A\mathcal{D}|D|^{-2-n} = -\sum_l a_{\alpha_2, -l} a_{\alpha_1, l} \text{Res}_{s=0} \left(\sum_k' \frac{k_{\mu_1} k_{\mu_2}}{|k|^{s+2+n}} \right) \text{Tr}(\gamma^{\alpha, \mu}).$$

We have also, with $\varepsilon JAJ^{-1} = iR(A_\alpha) \otimes \gamma^a$,

$$\text{Tr}(\varepsilon JAJ^{-1}\mathcal{D}A\mathcal{D}|D|^{-s-2-n}) \sim_c \sum_k' \sum_{l_1, l_2} a_{\alpha_2, l_2} a_{\alpha_1, l_1} \tau(U_{-k} U_{l_1} U_k U_{l_2}) \frac{k_{\mu_1} (k+l_1)_{\mu_2}}{|k|^{s+2+n}} \text{Tr}(\gamma^{\alpha, \mu}).$$

which gives

$$\int \varepsilon JAJ^{-1}\mathcal{D}A\mathcal{D}|D|^{-2-n} = a_{\alpha_2, 0} a_{\alpha_1, 0} \text{Res}_{s=0} \left(\sum_k' \frac{k_{\mu_1} k_{\mu_2}}{|k|^{s+2+n}} \right) \text{Tr}(\gamma^{\alpha, \mu}).$$

Thus,

$$\frac{1}{2} \int \tilde{A}\mathcal{D}\tilde{A}\mathcal{D}|D|^{-2-n} = \left(a_{\alpha_2, 0} a_{\alpha_1, 0} - \sum_l a_{\alpha_2, -l} a_{\alpha_1, l} \right) \text{Res}_{s=0} \left(\sum_k' \frac{k_{\mu_1} k_{\mu_2}}{|k|^{s+2+n}} \right) \text{Tr}(\gamma^{\alpha, \mu}).$$

With $\sum_k' \frac{k_{\mu_1} k_{\mu_2}}{|k|^{s+2+n}} = \frac{\delta_{\mu_1 \mu_2}}{n} Z_n(s+n)$ and $C_n := \text{Res}_{s=0} Z_n(s+n) = 2\pi^{n/2} \Gamma(n/2)^{-1}$ we obtain

$$\frac{1}{2} \int \tilde{A}\mathcal{D}\tilde{A}\mathcal{D}|D|^{-2-n} = \left(a_{\alpha_2, 0} a_{\alpha_1, 0} - \sum_l a_{\alpha_2, -l} a_{\alpha_1, l} \right) \frac{C_n}{n} \text{Tr}(\gamma^{\alpha_2} \gamma^{\mu} \gamma^{\alpha_1} \gamma^{\mu}).$$

Since $\text{Tr}(\gamma^{\alpha_2} \gamma^{\mu} \gamma^{\alpha_1} \gamma^{\mu}) = 2^m (2-n) \delta^{\alpha_2, \alpha_1}$, we get

$$\frac{1}{2} \int \tilde{A}\mathcal{D}\tilde{A}\mathcal{D}|D|^{-2-n} = 2^m \left(-a_{\alpha, 0} a_0^\alpha + \sum_l a_{\alpha, -l} a_l^\alpha \right) \frac{C_n (n-2)}{n}.$$

Equation (114) now proves the lemma. \square

Lemma 8.12. *If $\frac{1}{2\pi}\Theta$ is badly approximable, then for any $P \in \Psi_1(\mathcal{A})$ and $q \in \mathbb{N}$, q odd,*

$$\int P|D|^{-(n-q)} = 0.$$

Proof. There exist $B \in \mathcal{D}_1(\mathcal{A})$ and $p \in \mathbb{N}_0$ such that $P = BD^{-2p} + R$ where R is in OP^{-q-1} . As a consequence, $\int P|D|^{-(n-q)} = \int B|D|^{-n-2p+q}$. Assume $B = a_r b_r \mathcal{D}^{q_{r-1}} a_{r-1} b_{r-1} \cdots \mathcal{D}^{q_1} a_1 b_1$ where $r \in \mathbb{N}$, $a_i \in \mathcal{A}$, $b_i \in J\mathcal{A}J^{-1}$, $q_i \in \mathbb{N}$. If we prove that $\int B|D|^{-n-2p+q} = 0$, then the general case will follow by linearity. We note $a_i =: \sum_l a_{i,l} U_l$ and $b_i =: \sum_l b_{i,l} U_l$. With the shorthand $k_{\mu_1, \mu_{q_i}} := k_{\mu_1} \cdots k_{\mu_{q_i}}$ and $\gamma^{\mu_1, \mu_{q_i}} = \gamma^{\mu_1} \cdots \gamma^{\mu_{q_i}}$, we get

$$\mathcal{D}^{q_1} a_1 b_1 U_k \otimes e_j = \sum_{l_1, l'_1} a_{1, l_1} b_{1, l'_1} U_{l_1} U_k U_{l'_1} (k + l_1 + l'_1)_{\mu_1, \mu_{q_1}} \otimes \gamma^{\mu_1, \mu_{q_1}} e_j$$

which gives, after iteration,

$$B U_k \otimes e_j = \sum_{l, l'} \tilde{a}_l \tilde{b}_{l'} U_{l_r} \cdots U_{l_1} U_k U_{l'_1} \cdots U_{l'_r} \prod_{i=1}^{r-1} (k + \hat{l}_i + \hat{l}'_i)_{\mu_1^i, \mu_{q_i}^i} \otimes \gamma^{\mu_1^{r-1}, \mu_{q_{r-1}}^{r-1}} \cdots \gamma^{\mu_1^1, \mu_{q_1}^1} e_j$$

where $\tilde{a}_l := a_{1, l_1} \cdots a_{r, l_r}$ and $\tilde{b}_{l'} := b_{1, l'_1} \cdots b_{r, l'_r}$. Let's note $Q_\mu(k, l, l') := \prod_{i=1}^{r-1} (k + \hat{l}_i + \hat{l}'_i)_{\mu_1^i, \mu_{q_i}^i}$ and $\gamma^\mu := \gamma^{\mu_1^{r-1}, \mu_{q_{r-1}}^{r-1}} \cdots \gamma^{\mu_1^1, \mu_{q_1}^1}$. Thus,

$$\int B |D|^{-n-2p+q} = \text{Res}_{s=0} \sum_k' \sum_{l, l'} \tilde{a}_l \tilde{b}_{l'} \tau(U_{-k} U_{l_r} \cdots U_{l_1} U_k U_{l'_1} \cdots U_{l'_r}) \frac{Q_\mu(k, l, l')}{|k|^{s+2p+n-q}} \text{Tr}(\gamma^\mu).$$

Since $U_{l_r} \cdots U_{l_1} U_k = U_k U_{l_r} \cdots U_{l_1} e^{-i \sum_1^r l_i \cdot \Theta k}$, we get

$$\tau(U_{-k} U_{l_r} \cdots U_{l_1} U_k U_{l'_1} \cdots U_{l'_r}) = \delta_{\sum_1^r l_i + l'_i, 0} e^{i\phi(l, l')} e^{-i \sum_1^r l_i \cdot \Theta k}$$

where ϕ is a real valued function. Thus,

$$\begin{aligned} \int B |D|^{-n-2p+q} &= \text{Res}_{s=0} \sum_k' \sum_{l, l'} e^{i\phi(l, l')} \delta_{\sum_1^r l_i + l'_i, 0} \tilde{a}_l \tilde{b}_{l'} \frac{Q_\mu(k, l, l') e^{-i \sum_1^r l_i \cdot \Theta k}}{|k|^{s+2p+n-q}} \text{Tr}(\gamma^\mu) \\ &=: \text{Res}_{s=0} f_\mu(s) \text{Tr}(\gamma^\mu). \end{aligned}$$

We decompose $Q_\mu(k, l, l')$ as a sum $\sum_{h=0}^r M_{h, \mu}(l, l') Q_{h, \mu}(k)$ where $Q_{h, \mu}$ is a homogeneous polynomial in (k_1, \cdots, k_n) and $M_{h, \mu}(l, l')$ is a polynomial in $((l_1)_1, \cdots, (l_r)_n, (l'_1)_1, \cdots, (l'_r)_n)$.

Similarly, we decompose $f_\mu(s)$ as $\sum_{h=0}^r f_{h, \mu}(s)$. Theorem 7.5 (ii) entails that $f_{h, \mu}(s)$ extends meromorphically to the whole complex plane \mathbb{C} with only one possible pole for $s + 2p + n - q = n + d$ where $d := \deg Q_{h, \mu}$. In other words, if $d + q - 2p \neq 0$, $f_{h, \mu}(s)$ is holomorphic at $s = 0$. Suppose now $d + q - 2p = 0$ (note that this implies that d is odd, since q is odd by hypothesis), then, by Theorem 7.5 (ii)

$$\text{Res}_{s=0} f_{h, \mu}(s) = V \int_{u \in S^{n-1}} Q_{h, \mu}(u) dS(u)$$

where $V := \sum_{l, l' \in Z} M_{h, \mu}(l, l') e^{i\phi(l, l')} \delta_{\sum_1^r l_i + l'_i, 0} \tilde{a}_l \tilde{b}_{l'}$ and $Z := \{l, l' : \sum_{i=1}^r l_i = 0\}$. Since d is odd, $Q_{h, \mu}(-u) = -Q_{h, \mu}(u)$ and $\int_{u \in S^{n-1}} Q_{h, \mu}(u) dS(u) = 0$. Thus, $\text{Res}_{s=0} f_{h, \mu}(s) = 0$ in any case, which gives the result. \square

As we have seen, the crucial point of the preceding lemma is the decomposition of the numerator of the series $f_\mu(s)$ as polynomials in k . This has been possible because we restricted our pseudodifferential operators to $\Psi_1(\mathcal{A})$.

Proof of Proposition 8.6. The top element follows from Proposition 5.26 and according to (91),

$$\int |D|^{-n} = \operatorname{Res}_{s=0} \operatorname{Tr} (|D|^{-s-n}) = 2^m \operatorname{Res}_{s=0} Z_n(s+n) = \frac{2^{m+1}\pi^{n/2}}{\Gamma(n/2)}.$$

For the second equality, we get from Lemmas 8.8 and 5.23

$$\operatorname{Res}_{s=n-k} \zeta_{\mathcal{D}_A}(s) = \sum_{p=1}^k \sum_{r_1, \dots, r_p=0}^{k-p} h(n-k, r, p) \int \varepsilon^{r_1}(Y) \cdots \varepsilon^{r_p}(Y) |D|^{-(n-k)}.$$

Corollary 5.22 and Lemma 8.12 imply that $\int \varepsilon^{r_1}(Y) \cdots \varepsilon^{r_p}(Y) |D|^{-(n-k)} = 0$, which gives the result.

Last equality follows from Lemma 8.11 and Corollary 5.28. \square

8.4 The spectral action

Here is the main result of this section.

Theorem 8.13. *Consider the noncommutative torus $(C^\infty(\mathbb{T}_\Theta^n), \mathcal{H}, \mathcal{D})$ of dimension $n \in \mathbb{N}$ where $\frac{1}{2\pi}\Theta$ is a real $n \times n$ real skew-symmetric badly approximable matrix, and a selfadjoint one-form $A = L(-iA_\alpha) \otimes \gamma^\alpha$. Then, the full spectral action of $\mathcal{D}_A = \mathcal{D} + A + \epsilon JAJ^{-1}$ is*

$$\mathcal{S}(\mathcal{D}_A, f, \Lambda) = 4\pi f_2 \Lambda^2 + \mathcal{O}(\Lambda^{-2}),$$

(ii) for $n = 4$,

$$\mathcal{S}(\mathcal{D}_A, f, \Lambda) = 8\pi^2 f_4 \Lambda^4 - \frac{4\pi^2}{3} f(0) \tau(F_{\mu\nu} F^{\mu\nu}) + \mathcal{O}(\Lambda^{-2}),$$

(iii) More generally, in

$$\mathcal{S}(\mathcal{D}_A, f, \Lambda) = \sum_{k=0}^n f_{n-k} c_{n-k}(A) \Lambda^{n-k} + \mathcal{O}(\Lambda^{-1}),$$

$c_{n-2}(A) = 0$, $c_{n-k}(A) = 0$ for k odd. In particular, $c_0(A) = 0$ when n is odd.

This result (for $n = 4$) has also been obtained in [42] using the heat kernel method. It is however interesting to get the result via direct computations of (74) since it shows how this formula is efficient. As we will see, the computation of all the noncommutative integrals require a lot of technical steps. One of the main points, namely to isolate where the Diophantine condition on Θ is assumed, is outlined here.

Remark 8.14. Note that all terms must be gauge invariants, namely, according to (108), invariant by $A_\alpha \longrightarrow \gamma_u(A_\alpha) = uA_\alpha u^* + u\delta_\alpha(u^*)$. A particular case is $u = U_k$ where $U_k\delta_\alpha(U_k^*) = -ik_\alpha U_0$.

In the same way, note that there is no contradiction with the commutative case where, for any selfadjoint one-form A , $\mathcal{D}_A = \mathcal{D}$ (so A is equivalent to 0!), since we assume in Theorem 8.13 that Θ is badly approximable, so \mathcal{A} cannot be commutative.

Conjecture 8.15. The constant term of the spectral action of \mathcal{D}_A on the noncommutative n -torus is proportional to the constant term of the spectral action of $\mathcal{D}+A$ on the commutative n -torus.

Remark 8.16. The appearance of a Diophantine condition for Θ has been characterized in dimension 2 by Connes [23, Prop. 49] where in this case, $\Theta = \theta \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix}$ with $\theta \in \mathbb{R}$. In fact, the Hochschild cohomology $H(\mathcal{A}_\Theta, \mathcal{A}_\Theta^*)$ satisfies $\dim H^j(\mathcal{A}_\Theta, \mathcal{A}_\Theta^*) = 2$ (or 1) for $j = 1$ (or $j = 2$) if and only if the irrational number θ satisfies a Diophantine condition like $|1 - e^{i2\pi n\theta}|^{-1} = \mathcal{O}(n^k)$ for some k .

Recall that when the matrix Θ is quite irrational (the lattice generated by its columns is dense after translation by \mathbb{Z}^n , see [50, Def. 12.8]), then the C^* -algebra generated by \mathcal{A}_Θ is simple.

Remark 8.17. It is possible to generalize above theorem to the case $\mathcal{D} = -ig^\mu{}_\nu \delta_\mu \otimes \gamma^\nu$ instead of (103) when g is a positive definite constant matrix. The formulae in Theorem 8.13 are still valid, up to obvious modifications due to volume variation.

8.5 Computations of f

In order to get this theorem, let us prove a few technical lemmas.

We suppose from now on that Θ is a skew-symmetric matrix in $\mathcal{M}_n(\mathbb{R})$. No other hypothesis is assumed for Θ , except when it is explicitly stated.

When A is a selfadjoint one-form, we define for $n \in \mathbb{N}$, $q \in \mathbb{N}$, $2 \leq q \leq n$ and $\sigma \in \{-, +\}^q$

$$\begin{aligned} \mathbb{A}^+ &:= ADD^{-2}, \\ \mathbb{A}^- &:= \epsilon JAJ^{-1}DD^{-2}, \\ \mathbb{A}^\sigma &:= \mathbb{A}^{\sigma_q} \dots \mathbb{A}^{\sigma_1}. \end{aligned}$$

Lemma 8.18. We have for any $q \in \mathbb{N}$,

$$f(\tilde{A}D^{-1})^q = f(\tilde{A}DD^{-2})^q = \sum_{\sigma \in \{+, -\}^q} f \mathbb{A}^\sigma.$$

Proof. Since $P_0 \in OP^{-\infty}$, $D^{-1} = DD^{-2} \pmod{OP^{-\infty}}$ and $f(\tilde{A}D^{-1})^q = f(\tilde{A}DD^{-2})^q$. \square

Lemma 8.19. Let A be a selfadjoint one-form, $n \in \mathbb{N}$ and $q \in \mathbb{N}$ with $2 \leq q \leq n$ and $\sigma \in \{-, +\}^q$. Then

$$f \mathbb{A}^\sigma = f \mathbb{A}^{-\sigma}.$$

Definition 8.20. In [13] has been introduced the vanishing tadpole hypothesis:

$$\oint AD^{-1} = 0, \text{ for all } A \in \Omega_{\mathcal{D}}^1(\mathcal{A}). \quad (115)$$

By the following lemma, this condition is satisfied for the noncommutative torus.

Lemma 8.21. Let $n \in \mathbb{N}$, $A = L(-iA_{\alpha}) \otimes \gamma^{\alpha} = -i \sum_{l \in \mathbb{Z}^n} a_{\alpha,l} U_l \otimes \gamma^{\alpha}$, $A_{\alpha} \in \mathcal{A}_{\Theta}$, where $\{a_{\alpha,l}\}_{l \in \mathcal{S}(\mathbb{Z}^n)}$, be a hermitian one-form. Then,

(i) $\oint A^p D^{-q} = \oint (\epsilon JAJ^{-1})^p D^{-q} = 0$ for $p \geq 0$ and $1 \leq q < n$ (case $p = q = 1$ is tadpole hypothesis.)

(ii) If $\frac{1}{2\pi}\Theta$ is badly-approximable, then $\oint BD^{-q} = 0$ for $1 \leq q < n$ and any B in the algebra generated by \mathcal{A} , $[\mathcal{D}, \mathcal{A}]$, JAJ^{-1} and $J[\mathcal{D}, \mathcal{A}]J^{-1}$.

Proof. (i) Let us compute

$$\oint A^p (\epsilon JAJ^{-1})^{p'} D^{-q}.$$

With $A = L(-iA_{\alpha}) \otimes \gamma^{\alpha}$ and $\epsilon JAJ^{-1} = R(iA_{\alpha}) \otimes \gamma^{\alpha}$, we get

$$A^p = L(-iA_{\alpha_1}) \cdots L(-iA_{\alpha_p}) \otimes \gamma^{\alpha_1} \cdots \gamma^{\alpha_p}$$

and

$$(\epsilon JAJ^{-1})^{p'} = R(iA_{\alpha'_1}) \cdots R(iA_{\alpha'_{p'}}) \otimes \gamma^{\alpha'_1} \cdots \gamma^{\alpha'_{p'}}.$$

We note $\tilde{a}_{\alpha,l} := a_{\alpha_1,l_1} \cdots a_{\alpha_p,l_p}$. Since

$$L(-iA_{\alpha_1}) \cdots L(-iA_{\alpha_p}) R(iA_{\alpha'_1}) \cdots R(iA_{\alpha'_{p'}}) U_k = (-i)^p i^{p'} \sum_{l,l'} \tilde{a}_{\alpha,l} \tilde{a}_{\alpha',l'} U_{l_1} \cdots U_{l_p} U_k U_{l'_{p'}} \cdots U_{l'_1},$$

and

$$U_{l_1} \cdots U_{l_p} U_k = U_k U_{l_1} \cdots U_{l_p} e^{-i(\sum_i l_i) \cdot \Theta k},$$

we get, with

$$\begin{aligned} U_{l,l'} &:= U_{l_1} \cdots U_{l_p} U_{l'_{p'}} \cdots U_{l'_1}, \\ g_{\mu,\alpha,\alpha'}(s, k, l, l') &:= e^{ik \cdot \Theta \sum_j l_j} \frac{k_{\mu_1} \cdots k_{\mu_q}}{|k|^{s+2q}} \tilde{a}_{\alpha,l} \tilde{a}_{\alpha',l'}, \\ \gamma^{\alpha,\alpha',\mu} &:= \gamma^{\alpha_1} \cdots \gamma^{\alpha_p} \gamma^{\alpha'_1} \cdots \gamma^{\alpha'_{p'}} \gamma^{\mu_1} \cdots \gamma^{\mu_q}, \end{aligned}$$

$$A^p (\epsilon JAJ^{-1})^{p'} D^{-q} |D|^{-s} U_k \otimes e_i \sim_c (-i)^p i^{p'} \sum_{l,l'} g_{\mu,\alpha,\alpha'}(s, k, l, l') U_k U_{l,l'} \otimes \gamma^{\alpha,\alpha',\mu} e_i.$$

Thus, $\oint A^p (\epsilon JAJ^{-1})^{p'} D^{-q} = \text{Res}_{s=0} f(s)$ where

$$\begin{aligned} f(s) &:= \text{Tr} \left(A^p (\epsilon JAJ^{-1})^{p'} D^{-q} |D|^{-s} \right) \\ &\sim_c (-i)^p i^{p'} \sum'_{k \in \mathbb{Z}^n} \langle U_k \otimes e_i, \sum_{l,l'} g_{\mu,\alpha,\alpha'}(s, k, l, l') U_k U_{l,l'} \otimes \gamma^{\alpha,\alpha',\mu} e_i \rangle \\ &\sim_c (-i)^p i^{p'} \sum'_{k \in \mathbb{Z}^n} \tau \left(\sum_{l,l'} g_{\mu,\alpha,\alpha'}(s, k, l, l') U_{l,l'} \right) \text{Tr}(\gamma^{\mu,\alpha,\alpha'}) \\ &\sim_c (-i)^p i^{p'} \sum'_{k \in \mathbb{Z}^n} \sum_{l,l'} g_{\mu,\alpha,\alpha'}(s, k, l, l') \tau(U_{l,l'}) \text{Tr}(\gamma^{\mu,\alpha,\alpha'}). \end{aligned}$$

It is straightforward to check that the series $\sum'_{k,l,l'} g_{\mu,\alpha,\alpha'}(s,k,l,l') \tau(U_{l,l'})$ is absolutely summable if $\Re(s) > R$ for a $R > 0$. Thus, we can exchange the summation on k and l, l' , which gives

$$f(s) \sim_c (-i)^p i^{p'} \sum_{l,l'} \sum'_{k \in \mathbb{Z}^n} g_{\mu,\alpha,\alpha'}(s,k,l,l') \tau(U_{l,l'}) \text{Tr}(\gamma^{\mu,\alpha,\alpha'}).$$

If we suppose now that $p' = 0$, we see that,

$$f(s) \sim_c (-i)^p \sum_l \sum'_{k \in \mathbb{Z}^n} \frac{k_{\mu_1} \dots k_{\mu_q}}{|k|^{s+2q}} \tilde{a}_{\alpha,l} \delta_{\sum_{i=1}^p l_i, 0} \text{Tr}(\gamma^{\mu,\alpha,\alpha'})$$

which is, by Proposition 7.16, analytic at 0. In particular, for $p = q = 1$, we see that $f AD^{-1} = 0$, i.e. the vanishing tadpole hypothesis is satisfied. Similarly, if we suppose $p = 0$, we get

$$f(s) \sim_c (-i)^{p'} \sum_{l'} \sum'_{k \in \mathbb{Z}^n} \frac{k_{\mu_1} \dots k_{\mu_q}}{|k|^{s+2q}} \tilde{a}_{\alpha,l'} \delta_{\sum_{i=1}^{p'} l'_i, 0} \text{Tr}(\gamma^{\mu,\alpha,\alpha'})$$

which is holomorphic at 0.

(ii) Adapting the proof of Lemma 8.12 to our setting (taking $q_i = 0$, and adding gamma matrices components), we see that

$$\oint B D^{-q} = \text{Res}_{s=0} \sum_k \sum'_{l,l'} e^{i\phi(l,l')} \delta_{\sum_1^r l_i + l'_i, 0} \tilde{a}_{\alpha,l} \tilde{b}_{\beta,l'} \frac{k_{\mu_1} \dots k_{\mu_q} e^{-i \sum_1^r l_i \Theta_k}}{|k|^{s+2q}} \text{Tr}(\gamma^{(\mu,\alpha,\beta)})$$

where $\gamma^{(\mu,\alpha,\beta)}$ is a complicated product of gamma matrices. By Theorem 7.5 (ii), since we suppose here that $\frac{1}{2\pi}\Theta$ is badly approximable, this residue is 0. \square

8.5.1 Even dimensional case

Corollary 8.22. *Same hypothesis as in Lemma 8.21.*

(i) Case $n = 2$:

$$\oint A^q D^{-q} = -\delta_{q,2} 4\pi \tau(A_\alpha A^\alpha).$$

(ii) Case $n = 4$: with the shorthand $\delta_{\mu_1, \dots, \mu_4} := \delta_{\mu_1 \mu_2} \delta_{\mu_3 \mu_4} + \delta_{\mu_1 \mu_3} \delta_{\mu_2 \mu_4} + \delta_{\mu_1 \mu_4} \delta_{\mu_2 \mu_3}$,

$$\oint A^q D^{-q} = \delta_{q,4} \frac{\pi^2}{12} \tau(A_{\alpha_1} \dots A_{\alpha_4}) \text{Tr}(\gamma^{\alpha_1} \dots \gamma^{\alpha_4} \gamma^{\mu_1} \dots \gamma^{\mu_4}) \delta_{\mu_1, \dots, \mu_4}.$$

Proof. (i, ii) The same computation as in Lemma 8.21 (i) (with $p' = 0$, $p = q = n$) gives

$$\oint A^n D^{-n} = \text{Res}_{s=0} (-i)^n \left(\sum_{k \in \mathbb{Z}^n} \frac{k_{\mu_1} \dots k_{\mu_n}}{|k|^{s+2n}} \right) \tau \left(\sum_{l \in (\mathbb{Z}^n)^n} \tilde{a}_{\alpha,l} U_{l_1} \dots U_{l_n} \right) \text{Tr}(\gamma^{\alpha_1} \dots \gamma^{\alpha_n} \gamma^{\mu_1} \dots \gamma^{\mu_n})$$

and the result follows from Proposition 7.16. \square

We will use few notations:

For $n \in \mathbb{N}$, $q \geq 2$, $l := (l_1, \dots, l_{q-1}) \in (\mathbb{Z}^n)^{q-1}$, $\alpha := (\alpha_1, \dots, \alpha_q) \in \{1, \dots, n\}^q$, $k \in \mathbb{Z}^n \setminus \{0\}$, $\sigma \in \{-, +\}^q$, $(a_i)_{1 \leq i \leq n} \in (\mathcal{S}(\mathbb{Z}^n))^n$,

$$\begin{aligned} l_q &:= - \sum_{1 \leq j \leq q-1} l_j, \quad \lambda_\sigma := (-i)^q \prod_{j=1 \dots q} \sigma_j, \quad \tilde{a}_{\alpha, l} := a_{\alpha_1, l_1} \dots a_{\alpha_q, l_q}, \\ \phi_\sigma(k, l) &:= \sum_{1 \leq j \leq q-1} (\sigma_j - \sigma_q) k \cdot \Theta l_j + \sum_{2 \leq j \leq q-1} \sigma_j (l_1 + \dots + l_{j-1}) \cdot \Theta l_j, \\ g_\mu(s, k, l) &:= \frac{k_{\mu_1} (k+l_1)_{\mu_2} \dots (k+l_1+\dots+l_{q-1})_{\mu_q}}{|k|^{s+2} |k+l_1|^2 \dots |k+l_1+\dots+l_{q-1}|^2}, \end{aligned}$$

with the convention $\sum_{2 \leq j \leq q-1} = 0$ when $q = 2$, and $g_\mu(s, k, l) = 0$ whenever $\widehat{l}_i = -k$ for a $1 \leq i \leq q-1$.

Lemma 8.23. *Let $A = L(-iA_\alpha) \otimes \gamma^\alpha = -i \sum_{l \in \mathbb{Z}^n} a_{\alpha, l} U_l \otimes \gamma^\alpha$ where $A_\alpha = -A_\alpha^* \in \mathcal{A}_\Theta$ and $\{a_{\alpha, l}\}_l \in \mathcal{S}(\mathbb{Z}^n)$, with $n \in \mathbb{N}$, be a hermitian one-form, and let $2 \leq q \leq n$, $\sigma \in \{-, +\}^q$.*

Then, $f \mathbb{A}^\sigma = \operatorname{Res}_{s=0} f(s)$ where

$$f(s) := \sum_{l \in (\mathbb{Z}^n)^{q-1}} \sum'_{k \in \mathbb{Z}^n} \lambda_\sigma e^{\frac{i}{2} \phi_\sigma(k, l)} g_\mu(s, k, l) \tilde{a}_{\alpha, l} \operatorname{Tr}(\gamma^{\alpha_q} \gamma^{\mu_q} \dots \gamma^{\alpha_1} \gamma^{\mu_1}).$$

In the following, we will use the shorthand

$$c := \frac{4\pi^2}{3}.$$

Lemma 8.24. *Suppose $n = 4$. Then, with the same hypothesis of Lemma 8.23,*

- (i) $\frac{1}{2} \mathcal{f}(\mathbb{A}^+)^2 = \frac{1}{2} \mathcal{f}(\mathbb{A}^-)^2 = c \sum_{l \in \mathbb{Z}^4} a_{\alpha_1, l} a_{\alpha_2, -l} (l^{\alpha_1} l^{\alpha_2} - \delta^{\alpha_1 \alpha_2} |l|^2).$
- (ii) $-\frac{1}{3} \mathcal{f}(\mathbb{A}^+)^3 = -\frac{1}{3} \mathcal{f}(\mathbb{A}^-)^3 = 4c \sum_{l_i \in \mathbb{Z}^4} a_{\alpha_3, -l_1 - l_2} a_{l_2}^{\alpha_1} a_{\alpha_1, l_1} \sin \frac{l_1 \cdot \Theta l_2}{2} l_1^{\alpha_3}.$
- (iii) $\frac{1}{4} \mathcal{f}(\mathbb{A}^+)^4 = \frac{1}{4} \mathcal{f}(\mathbb{A}^-)^4 = 2c \sum_{l_i \in \mathbb{Z}^4} a_{\alpha_1, -l_1 - l_2 - l_3} a_{\alpha_2, l_3} a_{l_2}^{\alpha_1} a_{l_1}^{\alpha_2} \sin \frac{l_1 \cdot \Theta (l_2 + l_3)}{2} \sin \frac{l_2 \cdot \Theta l_3}{2}.$

(iv) *Suppose $\frac{1}{2\pi} \Theta$ badly approximable. Then the crossed terms in $f(\mathbb{A}^+ + \mathbb{A}^-)^q$ vanish: if C is the set of all $\sigma \in \{-, +\}^q$ with $2 \leq q \leq 4$, such that there exist i, j satisfying $\sigma_i \neq \sigma_j$, we have $\sum_{\sigma \in C} f \mathbb{A}^\sigma = 0$.*

Lemma 8.25. *Suppose $n = 4$ and $\frac{1}{2\pi} \Theta$ badly approximable. For any self-adjoint one-form A ,*

$$\zeta_{D_A}(0) - \zeta_D(0) = -c \tau(F_{\alpha_1, \alpha_2} F^{\alpha_1 \alpha_2}).$$

Proof. By (47) and Lemma 8.18 we get

$$\zeta_{D_A}(0) - \zeta_D(0) = \sum_{q=1}^n \frac{(-1)^q}{q} \sum_{\sigma \in \{+, -\}^q} \mathcal{f} \mathbb{A}^\sigma.$$

By Lemma 8.24 (iv), we see that the crossed terms all vanish. Thus, with Lemma 8.19, we get

$$\zeta_{D_A}(0) - \zeta_D(0) = 2 \sum_{q=1}^n \frac{(-1)^q}{q} \mathcal{f}(\mathbb{A}^+)^q. \quad (116)$$

By definition,

$$\begin{aligned} F_{\alpha_1\alpha_2} &= i \sum_k (a_{\alpha_2,k} k_{\alpha_1} - a_{\alpha_1,k} k_{\alpha_2}) U_k + \sum_{k,l} a_{\alpha_1,k} a_{\alpha_2,l} [U_k, U_l] \\ &= i \sum_k \left[(a_{\alpha_2,k} k_{\alpha_1} - a_{\alpha_1,k} k_{\alpha_2}) - 2 \sum_l a_{\alpha_1,k-l} a_{\alpha_2,l} \sin\left(\frac{k \cdot \Theta l}{2}\right) \right] U_k. \end{aligned}$$

Thus

$$\begin{aligned} \tau(F_{\alpha_1\alpha_2} F^{\alpha_1\alpha_2}) &= \sum_{\alpha_1, \alpha_2=1}^{2^m} \sum_{k \in \mathbb{Z}^4} \left[(a_{\alpha_2,k} k_{\alpha_1} - a_{\alpha_1,k} k_{\alpha_2}) - 2 \sum_{l' \in \mathbb{Z}^4} a_{\alpha_1,k-l'} a_{\alpha_2,l'} \sin\left(\frac{k \cdot \Theta l'}{2}\right) \right] \\ &\quad \left[(a_{\alpha_2,-k} k_{\alpha_1} - a_{\alpha_1,-k} k_{\alpha_2}) - 2 \sum_{l'' \in \mathbb{Z}^4} a_{\alpha_1,-k-l''} a_{\alpha_2,l''} \sin\left(\frac{k \cdot \Theta l''}{2}\right) \right]. \end{aligned}$$

One checks that the term in a^q of $\tau(F_{\alpha_1\alpha_2} F^{\alpha_1\alpha_2})$ corresponds to the term $f(\mathbb{A}^+)^q$ given by Lemma 8.24. For $q = 2$, this is

$$-2 \sum_{l \in \mathbb{Z}^4, \alpha_1, \alpha_2} a_{\alpha_1,l} a_{\alpha_2,-l} (l_{\alpha_1} l_{\alpha_2} - \delta_{\alpha_1\alpha_2} |l|^2).$$

For $q = 3$, we compute the crossed terms:

$$i \sum_{k,k',l} (a_{\alpha_2,k} k_{\alpha_1} - a_{\alpha_1,k} k_{\alpha_2}) a_{k'}^{\alpha_1} a_l^{\alpha_2} (U_k [U_{k'}, l] + [U_{k'}, U_l] U_k),$$

which gives the following a^3 -term in $\tau(F_{\alpha_1\alpha_2} F^{\alpha_1\alpha_2})$

$$-8 \sum_{l_i} a_{\alpha_3,-l_1-l_2} a_{l_2}^{\alpha_1} a_{\alpha_1,l_1} \sin \frac{l_1 \cdot \Theta l_2}{2} l_1^{\alpha_3}.$$

For $q = 4$, this is

$$-4 \sum_{l_i} a_{\alpha_1,-l_1-l_2-l_3} a_{\alpha_2,l_3} a_{l_2}^{\alpha_1} a_{l_1}^{\alpha_2} \sin \frac{l_1 \cdot \Theta (l_2+l_3)}{2} \sin \frac{l_2 \cdot \Theta l_3}{2}$$

which corresponds to the term $f(\mathbb{A}^+)^4$. We get finally,

$$\sum_{q=1}^n \frac{(-1)^q}{q} \mathcal{f}(\mathbb{A}^+)^q = -\frac{c}{2} \tau(F_{\alpha_1,\alpha_2} F^{\alpha_1\alpha_2}). \quad (117)$$

Equations (116) and (117) yield the result. \square

Lemma 8.26. *Suppose $n = 2$. Then, with the same hypothesis as in Lemma 8.23,*

$$(i) \quad \mathcal{f}(\mathbb{A}^+)^2 = \mathcal{f}(\mathbb{A}^-)^2 = 0.$$

(ii) *Suppose $\frac{1}{2\pi}\Theta$ badly approximable. Then*

$$\mathcal{f} \mathbb{A}^+ \mathbb{A}^- = \mathcal{f} \mathbb{A}^- \mathbb{A}^+ = 0.$$

Lemma 8.27. *Suppose $n = 2$ and $\frac{1}{2\pi}\Theta$ badly approximable. Then, for any self-adjoint one-form A ,*

$$\zeta_{D_A}(0) - \zeta_D(0) = 0.$$

Proof. As in Lemma 8.25, we use (47) and Lemma 8.18 so the result follows from Lemma 8.26. \square

8.5.2 Odd dimensional case

Lemma 8.28. *Suppose n odd and $\frac{1}{2\pi}\Theta$ badly approximable. Then for any self-adjoint one-form A and $\sigma \in \{-, +\}^q$ with $2 \leq q \leq n$,*

$$\int \mathbb{A}^\sigma = 0.$$

Proof. Since $\mathbb{A}^\sigma \in \Psi_1(\mathcal{A})$, Lemma 8.12 with $k = n$ gives the result. \square

Corollary 8.29. *With the same hypothesis of Lemma 8.28, $\zeta_{D_A}(0) - \zeta_D(0) = 0$.*

Proof. As in Lemma 8.25, we use (47) and Lemma 8.18 so the result follows from Lemma 8.28. \square

8.6 Proof of the main result

Proof of Theorem 8.13.. (i) By (74) and Proposition 8.6, we get

$$\mathcal{S}(\mathcal{D}_A, f, \Lambda) = 4f_2 \Lambda^2 + f(0) \zeta_{D_A}(0) + \mathcal{O}(\Lambda^{-2}),$$

where $f_2 = \frac{1}{2} \int_0^\infty f(t) dt$. By Lemma 8.27, $\zeta_{D_A}(0) - \zeta_D(0) = 0$ and from Proposition 8.5, $\zeta_D(0) = 0$, so we get the result.

(ii) Similarly, $\mathcal{S}(\mathcal{D}_A, f, \Lambda) = 8\pi^2 f_4 \Lambda^4 + f(0) \zeta_{D_A}(0) + \mathcal{O}(\Lambda^{-2})$ with $f_4 = \frac{1}{2} \int_0^\infty f(t) t dt$. Lemma 8.25 implies that $\zeta_{D_A}(0) - \zeta_D(0) = -c\tau(F_{\mu\nu}F^{\mu\nu})$ and Proposition 8.5 yields the equality $\zeta_{D_A}(0) = -c\tau(F_{\mu\nu}F^{\mu\nu})$ and the result.

(iii) is a direct consequence of (74), Propositions 8.5, 8.6, and Corollary 8.29. \square

8.7 Beyond Diophantine equation

This section is an attempt to understand what happens if Θ is ‘in between’ rational numbers and “Diophantine numbers”. Consider the simplest case: \mathbb{T}^2 with

$$\Theta = \theta \begin{pmatrix} 0 & -1 \\ 1 & 0 \end{pmatrix},$$

To proceed, we need some results from number theory [8]:

Definition 8.30. *Let $f : \mathbb{R}_{\geq 1} \rightarrow \mathbb{R}_{> 0}$ be a continuous function such that $x \rightarrow x^2 f(x)$ is non-increasing. Consider the set*

$$\mathcal{F}(f) := \left\{ \theta \in \mathbb{R} : |\theta q - p| < q f(q) \text{ for infinitely many rational numbers } \frac{p}{q} \right\}.$$

The elements of $\mathcal{F}(f)$ are termed f -approximable.

Note that we cannot expect the above estimate to be valid for all rational numbers $\frac{p}{q}$ since for all irrational numbers θ , the set of fractional values of $(\theta q)_{q \geq 1}$ is dense in $[0, 1]$.

Theorem 8.31. *There exists an uncountable family of real numbers $\theta/(2\pi)$ which are f -approximable but not cf -approximable for any $0 < c < 1$.*

See [8, Exercise 1.5] for a proof.

Let us choose

$$f(x) = (2\pi x)^{-1} e^{-2x},$$

and fix a constant $c < 1$. Let us pick a θ which is f -approximable, but not cf -approximable. Consider now $g(t) := \text{Tr} \left(a J b J^{-1} e^{-t \mathcal{D}^2} \right)$. It is shown in [42] that, by tuning $a, b \in \mathcal{A}_\theta$, it is possible to make the difference $g(t) - g(t)_{\text{Dioph}}$ (of $g(t)$ and its value if we suppose that θ is a Diophantine number) of arbitrary order in t .

This shows how subtle can be the computation of spectral action!

9 The non-compact case

9.1 The matter is not only technical

When a Riemannian spin manifold M is non-compact, the Dirac operator, which exists as a selfadjoint extension when M is (geodesically) complete has no more a compact resolvent: its spectrum is not discrete but is \mathbb{R} ([47, Theorem 7.2.1] and similar results for hyperbolic spaces [47, p. 106].)

To see what appends, let us consider for instance the flat space $M = \mathbb{R}^d$ and the Hilbert space $\mathcal{H} = L^2(\mathbb{R}^d)$. Then the operator $f(x)g(-i\nabla)$ is formally given on ψ, ϕ in appropriate domains by

$$\langle \psi, f(x)g(-i\nabla)\phi \rangle = \int_{\mathbb{R}^d} \bar{f}(x)\phi(x)(g\widehat{\psi})^\vee(x)dx.$$

For $k \in \mathbb{Z}^d$, let χ_k be the characteristic function of the unit cube in \mathbb{R}^d with center at k and define for $p, q > 0$

$$\ell^q(L^p(\mathbb{R}^d)) := \{f \mid \|f\|_{p,q} := \left(\sum_k \|f\chi_k\|_p^q\right)^{1/q} < \infty\}$$

where $\|g\|_p := \left(\int_{\mathbb{R}^d} |g(x)|^p dx\right)^{1/p}$ is the usual norm of $L^p(\mathbb{R}^d)$.

Theorem 9.1. *Birman–Solomjak.*

(i) If $f, g \in \ell^p(L^2(\mathbb{R}^d))$ for $1 \leq p \leq 2$, then $f(x)g(-i\nabla)$ is in the Schatten class \mathcal{L}^p and $\|f(x)g(-i\nabla)\|_p \leq c_p \|f\|_{2,p} \|g\|_{2,p}$.

(i) If f, g are non zero, then $f(x)g(-i\nabla) \in \mathcal{L}^1(\mathcal{H})$ if and only if f and g are in $\ell^1(\mathcal{H})$.

For a proof, see [102, Chapter 4].

This shows that even if $g(x) = e^{-tx^2}$, the heat kernel $e^{-t\Delta}$ is never trace-class since $f = 1$ is not in $\ell^1(L^2(\mathbb{R}^d))$.

Thus, to cover at least the non-compact manifold case, Definition 5.1 has to be improved:

Definition 9.2. A non-compact spectral triple $(\mathcal{A}, \mathcal{H}, \mathcal{D})$ is the data of an involutive algebra \mathcal{A} with a faithful representation π on a Hilbert space \mathcal{H} , a preferred unitization $\tilde{\mathcal{A}}$ of \mathcal{A} and a selfadjoint operator \mathcal{D} such that

- $a(\mathcal{D} - \lambda)^{-1}$ is compact for all $a \in \mathcal{A}$ and $\lambda \notin \text{Sp } \mathcal{D}$.
- $[\mathcal{D}, \pi(a)]$ is bounded for any $a \in \tilde{\mathcal{A}}$.

All definitions of regularity, finiteness and orientation have to be modified with $\tilde{\mathcal{A}}$ instead of \mathcal{A} , see also [9].

In the first constraint of this definition we recover a certain discreteness which, with $a = 1$, is the compact case (the algebra can have a unit). This matter is not only technical since now there is a deeper intertwining of the choice of the algebra \mathcal{A} and the operator \mathcal{D} to get a spectral triple. Moreover, a tentative of modification of \mathcal{D} is quite often forbidden by the second constraint.

The case of non-compact spin manifold has been considered by Rennie [91–93]. This has been improved in [40] which studied the Moyal plane. Actually, a compactification of this plane is the noncommutative torus!

9.2 The Moyal product

Reference: [40].

For any finite dimension k , let Θ be a real skewsymmetric $k \times k$ matrix, let $s \cdot t$ denote the usual scalar product on Euclidean \mathbb{R}^k and let $\mathcal{S}(\mathbb{R}^k)$ be the space of complex Schwartz functions on \mathbb{R}^k . One defines, for $f, h \in \mathcal{S}(\mathbb{R}^k)$, the corresponding Moyal or twisted product:

$$f \star_{\Theta} h(x) := (2\pi)^{-k} \iint f(x - \frac{1}{2}\Theta u) h(x + t) e^{-iu \cdot t} d^k u d^k t. \quad (118)$$

In Euclidean field theory, the entries of Θ have the dimensions of an area. Because Θ is skewsymmetric, complex conjugation reverses the product: $(f \star_{\Theta} h)^* = h^* \star_{\Theta} f^*$.

Assume Θ to be nondegenerate, that is to say, $\sigma(s, t) := s \cdot \Theta t$ to be symplectic. This implies even dimension, $k = 2N$. We note that Θ^{-1} is also skewsymmetric; let $\theta > 0$ be defined by $\theta^{2N} := \det \Theta$. Then formula (118) may be rewritten as

$$f \star_{\Theta} h(x) = (\pi\theta)^{-2N} \iint f(x + s) h(x + t) e^{-2is \cdot \Theta^{-1} t} d^{2N} s d^{2N} t. \quad (119)$$

This form is very familiar from phase-space quantum mechanics, where \mathbb{R}^{2N} is parametrized by N conjugate pairs of position and momentum variables, and the entries of Θ have the dimensions of an action; one then selects $\Theta = \hbar S := \hbar \begin{pmatrix} 0 & 1_N \\ -1_N & 0 \end{pmatrix}$. Indeed, the product \star (or rather, its commutator) was introduced in that context by Moyal [82], using a series development in powers of \hbar whose first nontrivial term gives the Poisson bracket; later, it was rewritten in the above integral form. These are actually oscillatory integrals, of which Moyal's series development,

$$f \star_{\hbar} g(x) = \sum_{\alpha \in \mathbb{N}^{2N}} \left(\frac{i\hbar}{2}\right)^{|\alpha|} \frac{1}{\alpha!} \frac{\partial f}{\partial x^{\alpha}}(x) \frac{\partial g}{\partial (Sx)^{\alpha}}(x), \quad (120)$$

is an asymptotic expansion. The first integral form (118) of the Moyal product was exploited by Rieffel in a remarkable monograph [96], who made it the starting point for a more general deformation theory of C^* -algebras.

With the choice $\Theta = \theta S$ made, the Moyal product can also be written

$$f \star_{\theta} g(x) := (\pi\theta)^{-2N} \iint f(y) g(z) e^{\frac{2i}{\theta}(x-y) \cdot S(x-z)} d^{2N} y d^{2N} z. \quad (121)$$

Of course, our definitions make sense only under certain hypotheses on f and g [49, 108].

Lemma 9.3. [49] *Let $f, g \in \mathcal{S}(\mathbb{R}^{2N})$. Then*

- (i) $f \star_{\theta} g \in \mathcal{S}(\mathbb{R}^{2N})$.
- (ii) \star_{θ} is a bilinear associative product on $\mathcal{S}(\mathbb{R}^{2N})$. Moreover, complex conjugation of functions $f \mapsto f^*$ is an involution for \star_{θ} .
- (iii) Let $j = 1, 2, \dots, 2N$. The Leibniz rule is satisfied:

$$\frac{\partial}{\partial x_j} (f \star_{\theta} g) = \frac{\partial f}{\partial x_j} \star_{\theta} g + f \star_{\theta} \frac{\partial g}{\partial x_j}. \quad (122)$$

(iv) Pointwise multiplication by any coordinate x_j obeys

$$x_j(f \star_\theta g) = f \star_\theta (x_j g) + \frac{i\theta}{2} \frac{\partial f}{\partial (Sx)_j} \star_\theta g = (x_j f) \star_\theta g - \frac{i\theta}{2} f \star_\theta \frac{\partial g}{\partial (Sx)_j}. \quad (123)$$

(v) The product has the tracial property:

$$\langle f, g \rangle := \frac{1}{(\pi\theta)^N} \int f \star_\theta g(x) d^{2N}x = \frac{1}{(\pi\theta)^N} \int g \star_\theta f(x) d^{2N}x = \frac{1}{(\pi\theta)^N} \int f(x) g(x) d^{2N}x.$$

(vi) Let $L_f^\theta \equiv L^\theta(f)$ be the left multiplication $g \mapsto f \star_\theta g$. Then $\lim_{\theta \downarrow 0} L_f^\theta g(x) = f(x)g(x)$, for $x \in \mathbb{R}^{2N}$.

Property (vi) is a consequence of the distributional identity

$$\lim_{\varepsilon \downarrow 0} \varepsilon^{-k} e^{ia \cdot b/\varepsilon} = (2\pi)^k \delta(a) \delta(b),$$

for $a, b \in \mathbb{R}^k$; convergence takes place in the standard topology [100] of $\mathcal{S}(\mathbb{R}^{2N})$. To simplify notation, we put $\mathcal{S} := \mathcal{S}(\mathbb{R}^{2N})$ and let $\mathcal{S}' := \mathcal{S}'(\mathbb{R}^{2N})$ be the dual space of tempered distributions. In view of (vi), we may denote by L_f^0 the pointwise product by f .

Theorem 9.4. [49] $\mathcal{A}_\theta := (\mathcal{S}, \star_\theta)$ is a nonunital associative, involutive Fréchet algebra with a jointly continuous product and a distinguished faithful trace.

Definition 9.5. The algebra \mathcal{A}_θ has a natural basis of eigenvectors f_{mn} of the harmonic oscillator, indexed by $m, n \in \mathbb{N}^N$. If

$$H_l := \frac{1}{2}(x_l^2 + x_{l+N}^2) \quad \text{for } l = 1, \dots, N \quad \text{and} \quad H := H_1 + H_2 + \dots + H_N,$$

then the f_{mn} diagonalize these harmonic oscillator Hamiltonians:

$$\begin{aligned} H_l \star_\theta f_{mn} &= \theta(m_l + \frac{1}{2}) f_{mn}, \\ f_{mn} \star_\theta H_l &= \theta(n_l + \frac{1}{2}) f_{mn}. \end{aligned} \quad (124)$$

They may be defined by

$$f_{mn} := \frac{1}{\sqrt{\theta^{|m|+|n|} m! n!}} (a^*)^m \star_\theta f_{00} \star_\theta a^n, \quad (125)$$

where f_{00} is the Gaussian function $f_{00}(x) := 2^N e^{-2H/\theta}$, and the annihilation and creation functions respectively are

$$a_l := \frac{1}{\sqrt{2}}(x_l + ix_{l+N}) \quad \text{and} \quad a_l^* := \frac{1}{\sqrt{2}}(x_l - ix_{l+N}). \quad (126)$$

One finds that $a^n := a_1^{n_1} \dots a_N^{n_N} = a_1^{\star_\theta n_1} \star_\theta \dots \star_\theta a_N^{\star_\theta n_N}$.

Proposition 9.6. [49, p. 877] The algebra $(\mathcal{S}, \star_\theta)$ has the (nonunique) factorization property: for all $h \in \mathcal{S}$ there exist $f, g \in \mathcal{S}$ such that $h = f \star_\theta g$.

Lemma 9.7. [49, 108] Let $f, g \in L^2(\mathbb{R}^{2N})$. Then

- (i) For $\theta \neq 0$, $f \star_\theta g$ lies in $L^2(\mathbb{R}^{2N})$. Moreover, $f \star_\theta g$ is uniformly continuous.
- (ii) \star_θ is a bilinear associative product on $L^2(\mathbb{R}^{2N})$. The complex conjugation of functions $f \mapsto f^*$ is an involution for \star_θ .
- (iii) The linear functional $f \mapsto \int f(x) dx$ on \mathcal{S} extends to $f \circ_0(\mathbb{R}^{2N}) := L^2(\mathbb{R}^{2N}) \star_\theta L^2(\mathbb{R}^{2N})$, and the product has the tracial property:

$$\langle f, g \rangle := (\pi\theta)^{-N} \int f \star_\theta g(x) d^{2N}x = (\pi\theta)^{-N} \int g \star_\theta f(x) d^{2N}x = (\pi\theta)^{-N} \int f(x) g(x) d^{2N}x.$$

- (iv) $\lim_{\theta \downarrow 0} L_f^\theta g(x) = f(x) g(x)$ almost everywhere on \mathbb{R}^{2N} .

Definition 9.8. Let $A_\theta := \{T \in \mathcal{S}' : T \star_\theta g \in L^2(\mathbb{R}^{2N}) \text{ for all } g \in L^2(\mathbb{R}^{2N})\}$, provided with the operator norm $\|L^\theta(T)\|_{\text{op}} := \sup\{\|T \star_\theta g\|_2 / \|g\|_2 : 0 \neq g \in L^2(\mathbb{R}^{2N})\}$.

Obviously $\mathcal{A}_\theta = \mathcal{S} \hookrightarrow A_\theta$. But \mathcal{A}_θ is not dense in A_θ .

Note that $\mathcal{G}_{00} \subset A_\theta$. This is clear from the following estimate.

Lemma 9.9. [49] If $f, g \in L^2(\mathbb{R}^{2N})$, then $f \star_\theta g \in L^2(\mathbb{R}^{2N})$ and $\|L_f^\theta\|_{\text{op}} \leq (2\pi\theta)^{-N/2} \|f\|_2$.

Proof. Expand $f = \sum_{m,n} c_{mn} \alpha_{mn}$ and $g = \sum_{m,n} d_{mn} \alpha_{mn}$ with respect to the orthonormal basis $\{\alpha_{nm}\} := (2\pi\theta)^{-N/2} \{f_{nm}\}$ of $L^2(\mathbb{R}^{2N})$. Then

$$\begin{aligned} \|f \star_\theta g\|_2^2 &= (2\pi\theta)^{-2N} \left\| \sum_{m,l} \left(\sum_n c_{mn} d_{nl} \right) f_{ml} \right\|_2^2 = (2\pi\theta)^{-N} \sum_{m,l} \left| \sum_n c_{mn} d_{nl} \right|^2 \\ &\leq (2\pi\theta)^{-N} \sum_{m,j} |c_{mj}|^2 \sum_{k,l} |d_{kl}|^2 = (2\pi\theta)^{-N} \|f\|_2^2 \|g\|_2^2, \end{aligned}$$

on applying the Cauchy–Schwarz inequality. \square

Proposition 9.10. [108] $(A_\theta, \|\cdot\|_{\text{op}})$ is a unital C^* -algebra of operators on $L^2(\mathbb{R}^{2N})$, isomorphic to $\mathcal{L}(L^2(\mathbb{R}^N))$ and including $L^2(\mathbb{R}^{2N})$. Moreover, there is a continuous injection of $*$ -algebras $\mathcal{A}_\theta \hookrightarrow A_\theta$, but \mathcal{A}_θ is not dense in A_θ .

Proposition 9.11. \mathcal{A}_θ is a (nonunital) Fréchet pre- C^* -algebra.

Proof. We adapt the argument for the commutative case in [50, p. 135]. To show that \mathcal{A}_θ is stable under the holomorphic functional calculus, we need only check that if $f \in \mathcal{A}_\theta$ and $1 + f$ is invertible in A_θ^0 with inverse $1 + g$, then the quasi-inverse g of f must lie in \mathcal{A}_θ . From $f + g + f \star_\theta g = 0$, we obtain $f \star_\theta f + g \star_\theta f + f \star_\theta g \star_\theta f = 0$, and it is enough to show that $f \star_\theta g \star_\theta f \in \mathcal{A}_\theta$, since the previous relation then implies $g \star_\theta f \in \mathcal{A}_\theta$, and then $g = -f - g \star_\theta f \in \mathcal{A}_\theta$ also.

Now, $A_\theta \subset \mathcal{G}_{-r,0}$ for any $r > N$ [108, p. 886]. Since $f \in \mathcal{G}_{s,p+r} \cap \mathcal{G}_{qt}$, for s, t arbitrary and p, q positive, we conclude that $f \star_\theta g \star_\theta f \in \mathcal{G}_{s,p+r} \star_\theta \mathcal{G}_{-r,0} \star_\theta \mathcal{G}_{qt} \subset \mathcal{G}_{st}$; as $\mathcal{S} = \bigcap_{s,t \in \mathbb{R}} \mathcal{G}_{st}$, the proof is complete. \square

Lemma 9.12. If $f \in \mathcal{S}$, then L_f^θ is a regularizing Ψ DO.

Proof. From (118), one at once sees that left Moyal multiplication by f is the pseudodifferential operator on \mathbb{R}^{2N} with symbol $f(x - \frac{\theta}{2}S\xi)$. Clearly L_f^θ extends to a continuous linear map from $C^\infty(\mathbb{R}^{2N})' \hookrightarrow \mathcal{S}'$ to $C^\infty(\mathbb{R}^{2N})$. The lemma also follows from the inequality

$$|\partial_x^\alpha \partial_\xi^\beta f(x - \frac{\theta}{2}S\xi)| \leq C_{K\alpha\beta}(1 + |\xi|^2)^{(d-|\beta|)/2},$$

valid for all $\alpha, \beta \in \mathbb{N}^{2N}$, any compact $K \subset \mathbb{R}^{2N}$, and any $d \in \mathbb{R}$, since $f \in \mathcal{S}$. \square

Remark 9.13. *Unlike for the case of a compact manifold, regularizing Ψ DOs are not necessarily compact operators!*

9.3 The preferred unitization of the Schwartz Moyal algebra

Definition 9.14. *Following Schwartz, we denote $\mathcal{B} := \mathcal{O}_0$, the space of smooth functions bounded together with all derivatives.*

A unitization of \mathcal{A}_θ is given by the algebra $\tilde{\mathcal{A}}_\theta := (\mathcal{B}, \star_\theta)$. The inclusion of \mathcal{A}_θ in \mathcal{B} is not dense, but this is not needed. $\tilde{\mathcal{A}}_\theta$ contains the constant functions and the plane waves, but no nonconstant polynomials and no imaginary-quadratic exponentials, such as $e^{iax_1x_2}$ in the case $N = 1$ (we will see later the pertinence of this).

Since \mathcal{B} is a unital $*$ -algebra with the Moyal product,

Proposition 9.15. *$\tilde{\mathcal{A}}_\theta$ is a unital Fréchet pre- C^* -algebra.*

An advantage of $\tilde{\mathcal{A}}_\theta$ is that the covering relation of the noncommutative plane to the NC torus is made transparent. To wit, the smooth noncommutative torus algebra $C^\infty(\mathbb{T}_\Theta^{2N})$ seen in Section 8.1 can be embedded in \mathcal{B} as periodic functions (with a fixed period parallelogram). This is in fact a Hopf algebra homomorphism: recall that $C^\infty(\mathbb{T}_\Theta^{2N})$ is a cotriangular Hopf algebra by exploiting the integral form (118) of (a periodic version of) the Moyal product.

We finally note the main reason for suitability of $\tilde{\mathcal{A}}_\theta$, namely, that each $[\mathcal{D}, L^\theta(f) \otimes 1_{2N}]$ lies in $\mathcal{A}_\theta \otimes M_{2N}(\mathbb{C})$, for $f \in \tilde{\mathcal{A}}_\theta$ and \mathcal{D} the Dirac operator on \mathbb{R}^{2N} .

9.4 The commutative case

When $\Theta = 0$ the Moyal product is the ordinary product.

Let \mathcal{A} be some appropriate subalgebra of $C^\infty(M)$ and \mathcal{D} be the Dirac operator, with k equal to the ordinary dimension of the spin manifold $M = \mathbb{R}^k$. Let \mathcal{H} be the space of square-integrable spinors. Then $[\mathcal{D}, f] = \mathcal{D}(f)$, just as in the unital case, and so the boundedness of $[D, \mathcal{A}]$ is unproblematic. In order to check whether $(\mathcal{A}, \mathcal{H}, \mathcal{D}, \chi)$ is a spectral triple, one first needs to determine whether products of the form $f(|\mathcal{D}| + \varepsilon)^{-k}$ are compact operators of Dixmier trace class, whose Dixmier trace is (a standard multiple of) $\int f(x) d^k x$. This compactness condition is guaranteed in the flat space case (taking $\mathcal{A} = \mathcal{S}(\mathbb{R}^k)$, say) by celebrated estimates in scattering theory [102].

The summability condition is a bit tougher. The Cesàro summability theory of [38] establishes that, for a positive pseudodifferential operator H of order d , acting on spinors, the spectral density asymptotically behaves as

$$d_H(x, x; \lambda') \sim \frac{2^{\lfloor k/2 \rfloor}}{d(2\pi)^k} \left(WRes H^{-k/d}(\lambda')^{(k-d)/d} + \dots \right),$$

in the Cesàro sense. (If the operator is not positive, one uses the “four parts” argument.) In our case, $H = a(|\mathcal{D}| + \varepsilon)^{-k}$ is pseudodifferential of order $-k$, so

$$d_H(x, x; \lambda') \sim -\frac{2^{\lfloor k/2 \rfloor} \Omega_k a(x)}{k(2\pi)^k} (\lambda'^{-2} + \dots),$$

as $\lambda' \rightarrow \infty$ in the Cesàro sense; here Ω_k is the hyper-area of the unit sphere in \mathbb{R}^k . We independently know that H is compact, so on integrating the spectral density over x and over $0 \leq \lambda' \leq \lambda$, we get that the number of eigenvalues of H less than λ is

$$N_H(\lambda) \sim \frac{2^{\lfloor k/2 \rfloor} \Omega_k \int a(x) d^k x}{k(2\pi)^k} \lambda^{-1} \quad \text{as } \lambda \rightarrow \infty.$$

This holds in the ordinary asymptotic sense, and not merely the Cesàro sense, by the “sandwich” argument used in the proof of [38, Cor. 4.1]. So finally,

$$\lambda_m(H) \sim \frac{2^{\lfloor k/2 \rfloor} \Omega_k \int a(x) d^k x}{k(2\pi)^k} m^{-1} \quad \text{as } m \rightarrow \infty, \quad (127)$$

and the Dixmier traceability of $a(|\mathcal{D}| + \varepsilon)^{-k}$, plus the value of its trace, follow at once.

The rest is a long but almost trivial verification. For instance, J is the charge conjugation operator on spinors; the algebra (\mathcal{B}, \star_0) is a suitable compactification; the domain \mathcal{H}^∞ consists of the smooth spinors; and so on. Thus, we get the following

Theorem 9.16. *The triple $(\mathcal{S}(\mathbb{R}^k), L^2(\mathbb{R}^k) \otimes \mathbb{C}^{2^{\lfloor k/2 \rfloor}}, \mathcal{D})$ on \mathbb{R}^k defines a noncompact commutative geometry of spectral dimension k .*

What about the non-flat case (of a spin manifold such that \mathcal{D} is selfadjoint)? Mainly because the previous Cesàro summability argument is purely local, everything carries over, if we choose for \mathcal{A} the algebra of smooth and compactly supported functions. Of course, in some contexts it may be useful to demand that M also has conic exits.

9.5 The Moyal plane

Let $\mathcal{A} = (\mathcal{S}(\mathbb{R}^{2N}), \star_\theta)$, with preferred unitization $\tilde{\mathcal{A}} := (\mathcal{B}(\mathbb{R}^{2N}), \star_\theta)$. The Hilbert space will be $\mathcal{H} := L^2(\mathbb{R}^{2N}) \otimes \mathbb{C}^{2^N}$ of ordinary square-integrable spinors. The representation of \mathcal{A} is given by $\pi^\theta: \mathcal{A} \rightarrow \mathcal{L}(\mathcal{H}) : f \mapsto L_f^\theta \otimes 1_{2^N}$, where L_f^θ acts on the “reduced” Hilbert space $\mathcal{H}_r := L^2(\mathbb{R}^{2N})$. In other words, if $a \in \mathcal{A}$ and $\Psi \in \mathcal{H}$, to obtain $\pi^\theta(a)\Psi$ we just left Moyal multiply Ψ by a componentwise.

This operator $\pi^\theta(f)$ is bounded, since it acts diagonally on \mathcal{H} and $\|L_f^\theta\| \leq (2\pi\theta)^{-N/2} \|f\|_2$ was proved in Lemma 9.9. Under this action, the elements of \mathcal{H} get the lofty name of *Moyal spinors*.

The selfadjoint Dirac operator is not “deformed”: it will be the ordinary Euclidean Dirac operator $\mathcal{D} := -i \gamma^\mu \partial_\mu$, where the hermitian matrices $\gamma^1, \dots, \gamma^{2N}$ satisfying $\{\gamma^\mu, \gamma^\nu\} = +2 \delta^{\mu\nu}$ irreducibly represent the Clifford algebra $\mathcal{C}\ell \mathbb{R}^{2N}$ associated to (\mathbb{R}^{2N}, η) , with η the standard Euclidean metric.

As a grading operator χ we take the usual chirality associated to the Clifford algebra:

$$\chi := \gamma_{2N+1} := 1_{\mathcal{H}_r} \otimes (-i)^N \gamma^1 \gamma^2 \dots \gamma^{2N}.$$

The notation γ_{2N+1} is a nod to physicists’ γ_5 . Thus $\chi^2 = (-1)^N (\gamma^1 \dots \gamma^{2N})^2 = (-1)^{2N} = 1$ and $\chi \gamma^\mu = -\gamma^\mu \chi$.

The real structure J is chosen to be the usual charge conjugation operator for spinors on \mathbb{R}^{2N} endowed with an Euclidean metric. Here, we only assume that $J^2 = \pm 1$ according to the “sign table” (41) and that

$$J(1_{\mathcal{H}_r} \otimes \gamma^\mu)J^{-1} = -1_{\mathcal{H}_r} \otimes \gamma^\mu$$

which guarantees the other requirements of the table. In general, in a given representation, it can be written as

$$J := CK, \tag{128}$$

where C denotes a suitable $2^N \times 2^N$ unitary matrix and K means complex conjugation. An important property of J is

$$J(L^\theta(f^*) \otimes 1_{2^N})J^{-1} = R^\theta(f) \otimes 1_{2^N}, \tag{129}$$

where $R^\theta(f) \equiv R_f^\theta$ is the right Moyal multiplication by f ; this follows from the antilinearity of J and the reversal of the twisted product under complex conjugation.

Lemma 9.3(iii) implies that $[\mathcal{D}, \pi^\theta(f)] = -iL^\theta(\partial_\mu f) \otimes \gamma^\mu =: \pi^\theta(\mathcal{D}(f))$ and this is bounded for $f \in \tilde{\mathcal{A}}_\theta = \mathcal{B}(\mathbb{R}^{2N})$ just as in the commutative case.

9.5.1 The compactness condition

In this subsection and the next, the main tools are techniques developed some time ago for scattering theory problems, as summarized in Simon’s booklet [102, Chap. 4]. We adopt the convention that $\mathcal{L}^\infty(\mathcal{H}) := \mathcal{K}(\mathcal{H})$, with $\|A\|_\infty := \|A\|_{\text{op}}$.

Let $g \in L^\infty(\mathbb{R}^{2N})$. We define the operator $g(-i\nabla)$ on \mathcal{H}_r as

$$g(-i\nabla)\psi := \mathcal{F}^{-1}(g\mathcal{F}\psi),$$

where \mathcal{F} is the ordinary Fourier transform. More in detail, for ψ in the correct domain,

$$g(-i\nabla)\psi(x) = (2\pi)^{-2N} \iint e^{i\xi \cdot (x-y)} g(\xi)\psi(y) d^{2N}\xi d^{2N}y.$$

The inequality $\|g(-i\nabla)\psi\|_2 = \|\mathcal{F}^{-1}g\mathcal{F}\psi\|_2 \leq \|g\|_\infty \|\psi\|_2$ entails that $\|g(-i\nabla)\|_\infty \leq \|g\|_\infty$.

Theorem 9.17. *Let $f \in \mathcal{A}$ and $\lambda \notin \text{sp } \mathcal{D}$. Then, if $R_{\mathcal{D}}(\lambda)$ is the resolvent operator of \mathcal{D} , then $\pi^\theta(f)R_{\mathcal{D}}(\lambda)$ is compact.*

Thanks to the first resolvent equation, $R_{\mathcal{D}}(\lambda) = R_{\mathcal{D}}(\lambda') + (\lambda' - \lambda)R_{\mathcal{D}}(\lambda)R_{\mathcal{D}}(\lambda')$, we may assume that $\lambda = i\mu$ with $\mu \in \mathbb{R}^*$. The theorem will follow from a series of lemmas interesting in themselves.

Lemma 9.18. *If $f \in \mathcal{S}$ and $0 \neq \mu \in \mathbb{R}$, then*

$$\pi^\theta(f)R_{\mathcal{D}}(i\mu) \in \mathcal{K}(\mathcal{H}) \iff \pi^\theta(f)|R_{\mathcal{D}}(i\mu)|^2 \in \mathcal{K}(\mathcal{H}).$$

Proof. We know that $L^\theta(f)^* = L^\theta(f^*)$. The “only if” part is obvious since $R_{\mathcal{D}}(i\mu)$ is a bounded normal operator. Conversely, if $\pi^\theta(f)|R_{\mathcal{D}}(i\mu)|^2$ is compact, then the operator $\pi^\theta(f)|R_{\mathcal{D}}(i\mu)|^2\pi^\theta(f^*)$ is compact. Since an operator T is compact if and only if TT^* is compact, the proof is complete. \square

The usefulness of this lemma stems from the diagonal nature of the action of the operator $\pi^\theta(f)|R_p(i\mu)|^2$ on $\mathcal{H} = \mathcal{H}_r \otimes \mathbb{C}^{2N}$; so in our arguments it is feasible to replace \mathcal{H} by \mathcal{H}_r , $\pi^\theta(f)$ by L_f^θ , and to use the scalar Laplacian $-\Delta := -\sum_{\mu=1}^{2N} \partial_\mu^2$ instead of the square of the Dirac operator \mathcal{D}^2 .

Lemma 9.19. *When $f, g \in \mathcal{H}_r$, $L_f^\theta g(-i\nabla)$ is a Hilbert–Schmidt operator such that, for all real θ ,*

$$\|L_f^\theta g(-i\nabla)\|_2 = \|L_f^0 g(-i\nabla)\|_2 = (2\pi)^{-N} \|f\|_2 \|g\|_2.$$

Proof. To prove that an operator A with integral kernel K_A is Hilbert–Schmidt, it suffices to check that $\int |K_A(x, y)|^2 dx dy$ is finite, and this will be equal to $\|A\|_2^2$ [102, Thm. 2.11]. So we compute $K_{L^\theta(f)g(-i\nabla)}$. In view of Lemma 9.12,

$$[L^\theta(f)g(-i\nabla)\psi](x) = \frac{1}{(2\pi)^{2N}} \iint f(x - \frac{\theta}{2}S\xi) g(\xi)\psi(y) e^{i\xi \cdot (x-y)} d^{2N}\xi d^{2N}y.$$

Thus

$$K_{L^\theta(f)g(-i\nabla)}(x, y) = \frac{1}{(2\pi)^{2N}} \int f(x - \frac{\theta}{2}S\xi) g(\xi) e^{i\xi \cdot (x-y)} d^{2N}\xi,$$

and $\int |K_{L^\theta(f)g(-i\nabla)}(x, y)|^2 dx dy$ is given by

$$\begin{aligned} & \frac{1}{(2\pi)^{4N}} \int \cdots \int \bar{f}(x - \frac{\theta}{2}S\xi) \bar{g}(\xi) f(x - \frac{\theta}{2}S\zeta) g(\zeta) e^{i(x-y) \cdot (\zeta - \xi)} d^{2N}x d^{2N}y d^{2N}\zeta d^{2N}\xi \\ &= \frac{1}{(2\pi)^{2N}} \iint |f(x - \frac{\theta}{2}S\xi)|^2 |g(\xi)|^2 d^{2N}x d^{2N}\xi = (2\pi)^{-2N} \|f\|_2^2 \|g\|_2^2 < \infty. \quad \square \end{aligned}$$

Remark 9.20. *As a consequence, we get*

$$\|\cdot\|_2\text{-}\lim_{\theta \rightarrow 0} L_f^\theta g(-i\nabla) = L_f^0 g(-i\nabla).$$

Lemma 9.21. *If $f \in \mathcal{H}_r$ and $g \in L^p(\mathbb{R}^{2N})$ with $2 \leq p < \infty$, then $L_f^\theta g(-i\nabla) \in \mathcal{L}^p(\mathcal{H}_r)$ and*

$$\|L_f^\theta g(-i\nabla)\|_p \leq (2\pi)^{-N(1/2+1/p)} \theta^{-N(1/2-1/p)} \|f\|_2 \|g\|_p.$$

Proof. The case $p = 2$ (with equality) is just the previous lemma. For $p = \infty$, we estimate $\|L_f^\theta g(-i\nabla)\|_\infty \leq (2\pi\theta)^{-N/2} \|f\|_2 \|g\|_\infty$: since $\|L_f^\theta g(-i\nabla)\|_\infty \leq \|L_f^\theta\|_\infty \|g(-i\nabla)\|_\infty$, this follows from Lemma 9.9 and a previous remark.

Now use complex interpolation for $2 < p < \infty$. For that, we first note that we may suppose $g \geq 0$: defining the function a with $|a| = 1$ and $g = a|g|$, we see that

$$\begin{aligned} \|L_f^\theta g(-i\nabla)\|_2^2 &= \text{Tr}(|L_f^\theta g(-i\nabla)|^2) = \text{Tr}(\bar{g}(-i\nabla) L_{f^*}^\theta L_f^\theta g(-i\nabla)) \\ &= \text{Tr}(|g|(-i\nabla) \bar{a}(-i\nabla) L_{f^*}^\theta L_f^\theta a(-i\nabla) |g|(-i\nabla)) \\ &= \text{Tr}(\bar{a}(-i\nabla) |g|(-i\nabla) L_{f^*}^\theta L_f^\theta |g|(-i\nabla) a(-i\nabla)) \\ &= \text{Tr}(|L_f^\theta |g|(-i\nabla)|^2) = \|L_f^\theta |g|(-i\nabla)\|_2^2, \end{aligned}$$

and

$$\begin{aligned} \|L_f^\theta g(-i\nabla)\|_\infty &= \|L_f^\theta a(-i\nabla) |g|(-i\nabla)\|_\infty = \|L_f^\theta |g|(-i\nabla) a(-i\nabla)\|_\infty \\ &\leq \|L_f^\theta |g|(-i\nabla)\|_\infty \|a(-i\nabla)\|_\infty = \|L_f^\theta |g|(-i\nabla)\|_\infty. \end{aligned}$$

Secondly, for any positive, bounded function g with compact support, we define the maps:

$$F_p : z \mapsto L_f^\theta g^{zp}(-i\nabla) : S = \{z \in \mathbb{C} \mid 0 \leq \Re z \leq \frac{1}{2}\} \rightarrow \mathcal{L}(\mathcal{H}_r).$$

For all $y \in \mathbb{R}$, $F_p(iy) = L_f^\theta g^{iyp}(-i\nabla) \in \mathcal{L}^\infty(\mathcal{H}_r)$ by Lemma 9.19 since g , being compactly supported, lies in \mathcal{H}_r . Moreover, $\|F_p(iy)\|_\infty \leq (2\pi\theta)^{-N/2}\|f\|_2$.

Also, by Lemma 9.19, $F_p(\frac{1}{2} + iy) \in \mathcal{L}^2(\mathcal{H}_r)$ and $\|F_p(\frac{1}{2} + iy)\|_2 = (2\pi)^{-N}\|f\|_2\|g^{p/2}\|_2$. Then complex interpolation (see [90, Chap. 9] and [102]) yields $F(z) \in \mathcal{L}^{1/\Re z}(\mathcal{H}_r)$, for all z in the strip S . Moreover,

$$\|F_p(z)\|_{1/\Re z} \leq \|F(0)\|_\infty^{1-2\Re z} \|F(\frac{1}{2})\|_2^{2\Re z} = \|f\|_2 (2\pi\theta)^{-\frac{N}{2}(1-2\Re z)} (2\pi)^{-2N\Re z} \|g^{p/2}\|_2^{2\Re z},$$

and applying this result at $z = 1/p$, we get for such g :

$$\|L_f^\theta g(-i\nabla)\|_p = \|F(1/p)\|_p \leq (2\pi)^{-N(1/2+1/p)} \theta^{-N(1/2-1/p)} \|f\|_2 \|g\|_p.$$

We finish by using the density of compactly supported bounded functions in $L^p(\mathbb{R}^{2N})$. \square

Lemma 9.22. *If $f \in \mathcal{S}$ and $0 \neq \mu \in \mathbb{R}$, then $\pi^\theta(f) |R_{\mathcal{D}}(i\mu)|^2 \in \mathcal{L}^p$ for $p > N$.*

Proof. We see that

$$\pi^\theta(f) |R_{\mathcal{D}}(i\mu)|^2 = (L_f^\theta \otimes 1_{2N}) (\mathcal{D} - i\mu)^{-1} (\mathcal{D} + i\mu)^{-1} = L_f^\theta (-\partial^\nu \partial_\nu + \mu^2)^{-1} \otimes 1_{2N}.$$

So this operator acts diagonally on $\mathcal{H}_r \otimes \mathbb{C}^{2N}$ and Lemma 9.21 implies that

$$\left\| L_f^\theta (-\partial^\nu \partial_\nu + \mu^2)^{-1} \right\|_p \leq (2\pi)^{-N(1/2+1/p)} \theta^{-N(1/2-1/p)} \|f\|_2 \left(\int \frac{d^{2N}\xi}{(\xi^\nu \xi_\nu + \mu^2)^p} \right)^{1/p},$$

which is finite for $p > N$. \square

Proof of Theorem 9.17. By Lemma 9.18, it was enough to prove that $\pi^\theta(f) |R_{\mathcal{D}}(i\mu)|^2$ is compact for a nonzero real μ . \square

9.5.2 Spectral dimension of the Moyal planes

Theorem 9.23. *The spectral dimension of the Moyal $2N$ -plane spectral triple is $2N$.*

We shall first establish existence properties.

Thanks to Lemma 9.21 and because $[\mathcal{D}, \pi^\theta(f)] = -iL^\theta(\partial_\mu f) \otimes \gamma^\mu$, we see that $\pi^\theta(f)(\mathcal{D}^2 + \varepsilon^2)^{-l}$ and $[\mathcal{D}, \pi^\theta(f)](\mathcal{D}^2 + \varepsilon^2)^{-l}$ lie in $\mathcal{L}^p(\mathcal{H})$ whenever $p > N/l$ (we always assume $\varepsilon > 0$). In the next lemma, we show that $[[\mathcal{D}], \pi^\theta(f)](\mathcal{D}^2 + \varepsilon^2)^{-l}$ has the same property of summability; this will become our main technical instrument for the subsection.

Lemma 9.24. *If $f \in \mathcal{S}$ and $\frac{1}{2} \leq l \leq N$, then $[[\mathcal{D}], \pi^\theta(f)](\mathcal{D}^2 + \varepsilon^2)^{-l} \in \mathcal{L}^p(\mathcal{H})$ for $p > N/l$.*

Proof. We use the following spectral identity for a positive operator A :

$$A = \frac{1}{\pi} \int_0^\infty \frac{A^2}{A^2 + \mu} \frac{d\mu}{\sqrt{\mu}},$$

and another identity for any operators A, B and $\lambda \notin \text{sp } A$:

$$[B, (A - \lambda)^{-1}] = (A - \lambda)^{-1}[A, B](A - \lambda)^{-1}. \quad (130)$$

Hence, for any $\rho > 0$,

$$\begin{aligned} [|\mathcal{D}|, \pi^\theta(f)] &= [|\mathcal{D}| + \rho, \pi^\theta(f)] = \frac{1}{\pi} \int_0^\infty \left[\frac{(|\mathcal{D}| + \rho)^2}{(|\mathcal{D}| + \rho)^2 + \mu}, \pi^\theta(f) \right] \frac{d\mu}{\sqrt{\mu}} \\ &= \frac{1}{\pi} \int_0^\infty \left(1 - \frac{(|\mathcal{D}| + \rho)^2}{(|\mathcal{D}| + \rho)^2 + \mu} \right) [(|\mathcal{D}| + \rho)^2, \pi^\theta(f)] \frac{1}{(|\mathcal{D}| + \rho)^2 + \mu} \frac{d\mu}{\sqrt{\mu}} \\ &= \frac{1}{\pi} \int_0^\infty \frac{1}{(|\mathcal{D}| + \rho)^2 + \mu} [(|\mathcal{D}| + \rho)^2, \pi^\theta(f)] \frac{1}{(|\mathcal{D}| + \rho)^2 + \mu} \sqrt{\mu} d\mu \\ &= \frac{1}{\pi} \int_0^\infty \frac{1}{(|\mathcal{D}| + \rho)^2 + \mu} \left(-\pi^\theta(\partial^\mu \partial_\mu f) - 2i(L^\theta(\partial_\mu f) \otimes \gamma^\mu)\mathcal{D} + 2\rho[|\mathcal{D}|, \pi^\theta(f)] \right) \\ &\quad \times \frac{1}{(|\mathcal{D}| + \rho)^2 + \mu} \sqrt{\mu} d\mu. \end{aligned} \quad (131)$$

This implies that

$$\begin{aligned} \left\| [|\mathcal{D}|, \pi^\theta(f)] (\mathcal{D}^2 + \varepsilon^2)^{-l} \right\|_p &\leq \frac{1}{\pi} \int_0^\infty \left\| \frac{1}{(|\mathcal{D}| + \rho)^2 + \mu} \left(-\pi^\theta(\partial^\mu \partial_\mu f) - 2i(L^\theta(\partial_\mu f) \otimes \gamma^\mu)\mathcal{D} \right. \right. \\ &\quad \left. \left. + 2\rho[|\mathcal{D}|, \pi^\theta(f)] \right) \frac{1}{(|\mathcal{D}| + \rho)^2 + \mu} (\mathcal{D}^2 + \varepsilon^2)^{-l} \right\|_p \sqrt{\mu} d\mu. \end{aligned}$$

Thus, the proof reduces to show that for any $f \in \mathcal{S}$,

$$\frac{1}{\pi} \int_0^\infty \left\| \frac{1}{(|\mathcal{D}| + \rho)^2 + \mu} \pi^\theta(f) \mathcal{D} \frac{1}{(|\mathcal{D}| + \rho)^2 + \mu} (\mathcal{D}^2 + \varepsilon^2)^{-l} \right\|_p \sqrt{\mu} d\mu < \infty. \quad (132)$$

Since the Schatten p -norm is a symmetric norm, and since, as in the proof of Theorem 9.17, only the reduced Hilbert space is affected, expression (132) is bounded by

$$\begin{aligned} &\frac{1}{\pi} \int_0^\infty \left\| \frac{1}{(|\mathcal{D}| + \rho)^2 + \mu} \right\|^{3/2} \left\| \frac{\mathcal{D}}{(\mathcal{D}^2 + \varepsilon^2)^{1/2}} \right\| \left\| \pi^\theta(f) \frac{1}{(\mathcal{D}^2 + \varepsilon^2)^{l-1/2}} \frac{1}{(|\mathcal{D}| + \rho)^2 + \mu} \right\|_p \sqrt{\mu} d\mu \\ &\leq \frac{1}{\pi} \int_0^\infty \left\| \pi^\theta(f) (\mathcal{D}^2 + \varepsilon^2)^{-l+1/2} (|\mathcal{D}| + \rho)^2 + \mu)^{-1/2} \right\|_p \frac{\sqrt{\mu} d\mu}{(\mu + \rho^2)^{3/2}}. \end{aligned}$$

Thanks to Lemma 9.21, we can estimate the μ -dependence of the last p -norm:

$$\begin{aligned} &\left\| \pi^\theta(f) (|\mathcal{D}| + \rho)^2 + \mu)^{-1/2} (\mathcal{D}^2 + \varepsilon^2)^{-l+1/2} \right\|_p \\ &\leq (2\pi)^{-N(1/2+1/p)} \theta^{-N(1/2-1/p)} \|f\|_2 \left\| (|\xi| + \rho)^2 + \mu)^{-1/2} (|\xi|^2 + \varepsilon^2)^{-l+1/2} \right\|_p \\ &\leq C(p, \theta) \left\| (|\xi| + \rho)^2 + \mu)^{-1/2} \right\|_q \left\| (|\xi|^2 + \varepsilon^2)^{-l+1/2} \right\|_r; \end{aligned}$$

with $p^{-1} = q^{-1} + r^{-1}$ appropriately chosen, these integrals are finite for all $q > 2N$ and $r > 2N/(2l - 1)$; for $l = \frac{1}{2}$, take $r = \infty$ and $q = p$. For such values,

$$\begin{aligned} & \left\| \pi^\theta(f) ((|\mathcal{D}| + \rho)^2 + \mu)^{-1/2} (\mathcal{D}^2 + \varepsilon^2)^{-l+1/2} \right\|_p \\ & \leq C(p, \theta, N; f) \| (|\xi|^2 + \varepsilon^2)^{-l+1/2} \|_r \Omega_{2N}^{1/q} \left(\int_0^\infty \frac{R^{2N-1}}{((R+\rho)^2 + \mu)^{q/2}} dR \right)^{1/q} \\ & = C(p, \theta, N; f) \| (|\xi|^2 + \varepsilon^2)^{-l+1/2} \|_r \pi^{N/q} \frac{\Gamma^{1/q}(\frac{q}{2}-N)}{\Gamma^{1/q}(\frac{q}{2})} \mu^{-1/2+N/q} =: C'(p, q, \theta, N; f) \mu^{-1/2+N/q}. \end{aligned}$$

Finally, the integral (132) is less than

$$C'(p, q, \theta, N; f) \int_0^\infty \frac{\mu^{N/q}}{(\mu + \rho^2)^{3/2}} d\mu,$$

which is finite for $q > 2N$ and $p > N/l$. This concludes the proof. \square

Lemma 9.25. *If $f \in \mathcal{S}$, then $\pi^\theta(f) (|\mathcal{D}| + \varepsilon)^{-1} \pi^\theta(f^*) \in \mathcal{L}^{2N+}(\mathcal{H})$.*

Proof. This is an extension to the Moyal context of the renowned inequality by Cwikel [102]. As remarked before, it is possible to replace \mathcal{D}^2 by $-\Delta$, $\pi^\theta(f)$ by L_f^θ and \mathcal{H} by \mathcal{H}_r . Consider $g(-i\nabla) := (\sqrt{-\Delta} + \varepsilon)^{-1}$. Since g is positive, it can be decomposed as $g = \sum_{n \in \mathbb{Z}} g_n$ where

$$g_n(x) := \begin{cases} g(x) & \text{if } 2^{n-1} < g(x) \leq 2^n, \\ 0 & \text{otherwise.} \end{cases}$$

For each $n \in \mathbb{Z}$, let A_n and B_n be the two operators

$$A_n := \sum_{k \leq n} L_f^\theta g_k(-i\nabla) L_{f^*}^\theta, \quad B_n := \sum_{k > n} L_f^\theta g_k(-i\nabla) L_{f^*}^\theta.$$

We estimate the uniform norm of the first part:

$$\begin{aligned} \|A_n\|_\infty & \leq \|L_f^\theta\|^2 \left\| \sum_{k \leq n} g_k(-i\nabla) \right\|_\infty \leq (2\pi\theta)^{-N} \|f\|_2^2 \left\| \sum_{k \leq n} g_k \right\|_\infty \\ & \leq (2\pi\theta)^{-N} \|f\|_2^2 2^n =: 2^n c_1(\theta, N; f). \end{aligned}$$

The trace norm of B_n can be computed using Lemma 9.19:

$$\begin{aligned} \|B_n\|_1 & = \left\| \left(\sum_{k > n} g_k(-i\nabla) \right)^{1/2} L_{f^*}^\theta \right\|_2^2 = \left\| L_f^\theta \left(\sum_{k > n} g_k(-i\nabla) \right)^{1/2} \right\|_2^2 = (2\pi)^{-2N} \|f\|_2^2 \left\| \left(\sum_{k > n} g_k \right)^{1/2} \right\|_2^2 \\ & = (2\pi)^{-2N} \|f\|_2^2 \left\| \sum_{k > n} g_k \right\|_1 = (2\pi)^{-2N} \|f\|_2^2 \sum_{k > n} \|g_k\|_1 \\ & \leq (2\pi)^{-2N} \|f\|_2^2 \sum_{k > n} \|g_k\|_\infty \nu\{\text{supp}(g_k)\}, \end{aligned}$$

where ν is the Lebesgue measure on \mathbb{R}^{2N} . By definition, $\|g_k\|_\infty \leq 2^k$ and

$$\begin{aligned} \nu\{\text{supp}(g_k)\} & = \nu\{x \in \mathbb{R}^{2N} : 2^{k-1} < g(x) \leq 2^k\} \leq \nu\{x \in \mathbb{R}^{2N} \mid (|x| + \varepsilon)^{-1} \geq 2^{k-1}\} \\ & \leq 2^{2N(1-k)} c_2. \end{aligned}$$

Therefore

$$\begin{aligned}\|B_n\|_1 &\leq (2\pi)^{-2N} \|f\|_2^2 2^{2N} c_2 \sum_{k>n} 2^{k(1-2N)} \\ &< \pi^{-2N} c_2 \|f\|_2^2 2^{n(1-2N)} =: 2^{n(1-2N)} c_3(N; f),\end{aligned}$$

where the second inequality follows because $N > \frac{1}{2}$.

We can now estimate the m th singular value μ_m of B_n (arranged in decreasing order with multiplicity): $\|B_n\|_1 = \sum_{k=0}^{\infty} \mu_k(B_n)$. Note that, for $m = 1, 2, 3, \dots$, we get that $\|B_n\|_1 \geq \sum_{k=0}^{m-1} \mu_k(B_n) \geq m \mu_m(B_n)$. Thus, $\mu_m(B_n) \leq \|B_n\|_1 m^{-1} \leq 2^{n(1-2N)} c_3 m^{-1}$. Now Fan's inequality [102, Thm. 1.7] yields

$$\begin{aligned}\mu_m(L_f^\theta g(-i\nabla) L_{f^*}^\theta) &= \mu_m(A_n + B_n) \leq \mu_1(A_n) + \mu_m(B_n) \\ &\leq \|A_n\| + \|B_n\|_1 m^{-1} \leq 2^n c_1 + 2^{n(1-2N)} c_3 m^{-1}.\end{aligned}$$

Given m , choose $n \in \mathbb{Z}$ so that $2^n \leq m^{-1/2N} < 2^{n+1}$. Then

$$\mu_m(L_f^\theta g(-i\nabla) L_{f^*}^\theta) \leq c_1 m^{-1/2N} + c_3 m^{-(1-2N)/2N} m^{-1} =: c_4(\theta, N; f) m^{-1/2N}.$$

Therefore $L_f^\theta (\sqrt{-\Delta} + \varepsilon)^{-1} L_{f^*}^\theta \in \mathcal{L}^{2N^+}(\mathcal{H}_r)$, and the statement of the lemma follows. \square

Corollary 9.26. *If $f, g \in \mathcal{S}$, then $\pi^\theta(f) (|\mathcal{D}| + \varepsilon)^{-1} \pi^\theta(g) \in \mathcal{L}^{2N^+}(\mathcal{H})$.*

Proof. Consider $\pi^\theta(f \pm g^*) (|\mathcal{D}| + \varepsilon)^{-1} \pi^\theta(f^* \pm g)$ and $\pi^\theta(f \pm ig^*) (|\mathcal{D}| + \varepsilon)^{-1} \pi^\theta(f^* \mp ig)$. \square

Corollary 9.27. *If $h \in \mathcal{S}$, then $\pi^\theta(h) (|\mathcal{D}| + \varepsilon)^{-1} \in \mathcal{L}^{2N^+}(\mathcal{H})$.*

Proof. Let $h = f \star_\theta g$. Then

$$\pi^\theta(h) (|\mathcal{D}| + \varepsilon)^{-1} = \pi^\theta(f) (|\mathcal{D}| + \varepsilon)^{-1} \pi^\theta(g) + \pi^\theta(f) [\pi^\theta(g), (|\mathcal{D}| + \varepsilon)^{-1}],$$

and we obtain from the identity (130) that

$$\pi^\theta(h) (|\mathcal{D}| + \varepsilon)^{-1} = \pi^\theta(f) (|\mathcal{D}| + \varepsilon)^{-1} \pi^\theta(g) + \pi^\theta(f) (|\mathcal{D}| + \varepsilon)^{-1} [|\mathcal{D}|, \pi^\theta(g)] (|\mathcal{D}| + \varepsilon)^{-1}.$$

By arguments similar to those of lemmata 9.21 and 9.24, the last term belongs to \mathcal{L}^p for $p > N$, and thus to \mathcal{L}^{2N^+} . \square

Boundedness of $(|\mathcal{D}| + \varepsilon)(\mathcal{D}^2 + \varepsilon^2)^{-1/2}$ follows from elementary Fourier analysis. And so the last corollary means that the spectral triple is “ $2N^+$ -summable”. We have taken care of the first assertion of the theorem. The next lemma is the last property of existence that we need.

Lemma 9.28. *If $f \in \mathcal{S}$, then $\pi^\theta(f)(|\mathcal{D}| + \varepsilon)^{-2N}$ and $\pi^\theta(f)(\mathcal{D}^2 + \varepsilon^2)^{-N}$ are in $\mathcal{L}^{1^+}(\mathcal{H})$.*

Proof. It suffices to prove that $\pi^\theta(f)(|\mathcal{D}| + \varepsilon)^{-2N} \in \mathcal{L}^{1^+}(\mathcal{H})$. We factorize $f \in \mathcal{S}$ according to Proposition 9.6, with the following notation:

$$\begin{aligned}f &= f_1 \star_\theta f_2 = f_1 \star_\theta f_{21} \star_\theta f_{22} = f_1 \star_\theta f_{21} \star_\theta f_{221} \star_\theta f_{222} \\ &= \dots = f_1 \star_\theta f_{21} \star_\theta f_{221} \star_\theta \dots \star_\theta f_{22\dots 21} \star_\theta f_{22\dots 22}.\end{aligned}$$

Therefore,

$$\begin{aligned} \pi^\theta(f) (|\mathcal{D}| + \varepsilon)^{-2N} &= \pi^\theta(f_1) (|\mathcal{D}| + \varepsilon)^{-1} \pi^\theta(f_2) (|\mathcal{D}| + \varepsilon)^{-2N+1} \\ &\quad + \pi^\theta(f_1) (|\mathcal{D}| + \varepsilon)^{-1} [|\mathcal{D}|, \pi^\theta(f_2)] (|\mathcal{D}| + \varepsilon)^{-2N}. \end{aligned} \quad (133)$$

By Lemma 9.21, $\pi^\theta(f_1)(|\mathcal{D}| + \varepsilon)^{-1} \in \mathcal{L}^p(\mathcal{H})$ whenever $p > 2N$; and by Lemma 9.24, the term $[|\mathcal{D}|, \pi^\theta(f_2)](|\mathcal{D}| + \varepsilon)^{-2N}$ lies in $\mathcal{L}^q(\mathcal{H})$ for $q > 1$. Hence, the last term on the right hand side of equation (133) lies in $\mathcal{L}^1(\mathcal{H})$. We may write the following equivalence relation:

$$\pi^\theta(f)(|\mathcal{D}| + \varepsilon)^{-2N} \sim \pi^\theta(f_1)(|\mathcal{D}| + \varepsilon)^{-1} \pi^\theta(f_2)(|\mathcal{D}| + \varepsilon)^{-2N+1},$$

where $A \sim B$ for $A, B \in \mathcal{K}(\mathcal{H})$ means that $A - B$ is trace-class. Thus,

$$\begin{aligned} \pi^\theta(f)(|\mathcal{D}| + \varepsilon)^{-2N} &\sim \pi^\theta(f_1)(|\mathcal{D}| + \varepsilon)^{-1} \pi^\theta(f_2)(|\mathcal{D}| + \varepsilon)^{-2N+1} \\ &= \pi^\theta(f_1)(|\mathcal{D}| + \varepsilon)^{-1} \pi^\theta(f_{21})(|\mathcal{D}| + \varepsilon)^{-1} \pi^\theta(f_{22})(|\mathcal{D}| + \varepsilon)^{-2N+2} \\ &\quad + \pi^\theta(f_1)(|\mathcal{D}| + \varepsilon)^{-1} \pi^\theta(f_{21})(|\mathcal{D}| + \varepsilon)^{-1} [|\mathcal{D}|, \pi^\theta(f_{22})] (|\mathcal{D}| + \varepsilon)^{-2N+1} \\ &\sim \pi^\theta(f_1)(|\mathcal{D}| + \varepsilon)^{-1} \pi^\theta(f_{21})(|\mathcal{D}| + \varepsilon)^{-1} \pi^\theta(f_{22})(|\mathcal{D}| + \varepsilon)^{-2N+2} \sim \dots \\ &\sim \pi^\theta(f_1)(|\mathcal{D}| + \varepsilon)^{-1} \pi^\theta(f_{21})(|\mathcal{D}| + \varepsilon)^{-1} \pi^\theta(f_{221})(|\mathcal{D}| + \varepsilon)^{-1} \dots \pi^\theta(f_{22\dots 22})(|\mathcal{D}| + \varepsilon)^{-1}. \end{aligned}$$

The second equivalence relation holds because $\pi^\theta(f_1)(|\mathcal{D}| + \varepsilon)^{-1} \pi^\theta(f_{21})(|\mathcal{D}| + \varepsilon)^{-1} \in \mathcal{L}^p(\mathcal{H})$ for $p > N$ by Lemma 9.21, and $[|\mathcal{D}|, \pi^\theta(f_{22})](|\mathcal{D}| + \varepsilon)^{-2N+1} \in \mathcal{L}^q(\mathcal{H})$ for $q > 2N/(2N - 1)$ by Lemma 9.24 again. The other equivalences come from similar arguments. Corollary 9.26, the Hölder inequality (see [50, Prop. 7.16]) and the inclusion $\mathcal{L}^1(\mathcal{H}) \subset \mathcal{L}^{1+}(\mathcal{H})$ finally yield the result. \square

Now we go for the computation of the Dixmier trace. Using the regularized trace for a Ψ DO:

$$\mathrm{Tr}_\Lambda(A) := (2\pi)^{-2N} \iint_{|\xi| \leq \Lambda} \sigma[A](x, \xi) d^{2N} \xi d^{2N} x,$$

the result can be conjectured because $\lim_{\Lambda \rightarrow \infty} \mathrm{Tr}_\Lambda(\cdot) / \log(\Lambda^{2N})$ is heuristically linked with the Dixmier trace, and the following computation:

$$\begin{aligned} &\lim_{\Lambda \rightarrow \infty} \frac{1}{2N \log \Lambda} \mathrm{Tr}_\Lambda \left(\pi^\theta(f) (\mathcal{D}^2 + \varepsilon^2)^{-N} \right) \\ &= \lim_{\Lambda \rightarrow \infty} t \frac{2^N}{2N(2\pi)^{2N} \log \Lambda} \iint_{|\xi| \leq \Lambda} f(x - \frac{\theta}{2} S\xi) (|\xi|^2 + \varepsilon^2)^{-N} d^{2N} \xi d^{2N} x \\ &= \frac{2^N \Omega_{2N}}{2N(2\pi)^{2N}} \int f(x) d^{2N} x. \end{aligned}$$

This is precisely the same result of (127), in the commutative case, for $k = 2N$. However, to establish it rigorously in the Moyal context requires a subtler strategy. We shall compute the Dixmier trace of $\pi^\theta(f) (\mathcal{D}^2 + \varepsilon^2)^{-N}$ as the residue of the ordinary trace of a related meromorphic family of operators. In turn we are allowed to introduce the explicit symbol formula that will establish measurability [25, 50], too.

We seek first to verify that \mathcal{A}_θ has *analytical dimension* equal to $2N$; that is, for $f \in \mathcal{A}_\theta$ the operator $\pi^\theta(f) (\mathcal{D}^2 + \varepsilon^2)^{-z/2}$ is trace-class if $\Re z > 2N$.

Lemma 9.29. *If $f \in \mathcal{S}$, then $L_f^\theta (\mathcal{D}^2 + \varepsilon^2)^{-z/2}$ is trace-class for $\Re z > 2N$, and*

$$\mathrm{Tr}[L_f^\theta (\mathcal{D}^2 + \varepsilon^2)^{-z/2}] = (2\pi)^{-2N} \iint f(x) (|\xi|^2 + \varepsilon^2)^{-z/2} d^{2N}\xi d^{2N}x.$$

Proof. If $a(x, \xi) \in \mathcal{K}_p(\mathbb{R}^{2k})$, for $p < -k$, is the symbol of a pseudodifferential operator A , then the operator is trace-class and moreover

$$\mathrm{Tr} A = (2\pi)^{-k} \iint a(x, \xi) d^k x d^k \xi.$$

This is easily proved by taking $a \in \mathcal{S}(\mathbb{R}^{2k})$ first and extending the resulting formula by continuity.

In our case, the symbol formula for a product of Ψ DOs yields, for $p > N$,

$$\begin{aligned} \sigma[L_f^\theta (-\Delta + \varepsilon^2)^{-p}](x, \xi) &= \sum_{\alpha \in \mathbb{N}^N} \frac{(-i)^{|\alpha|}}{\alpha!} \partial_\xi^\alpha \sigma[L_f^\theta](x, \xi) \partial_x^\alpha \sigma[(-\Delta + \varepsilon^2)^{-p}](x, \xi) \\ &= \sigma[L_f^\theta](x, \xi) \sigma[(-\Delta + \varepsilon^2)^{-p}](x, \xi) \\ &= f(x - \frac{\theta}{2} S\xi) (|\xi|^2 + \varepsilon^2)^{-p}. \end{aligned}$$

Therefore, for $p > N$,

$$\begin{aligned} \mathrm{Tr}(L_f^\theta (-\Delta + \varepsilon^2)^{-p}) &= (2\pi)^{-2N} \iint f(x - \frac{\theta}{2} S\xi) (|\xi|^2 + \varepsilon^2)^{-p} d^{2N}\xi d^{2N}x \\ &= (2\pi)^{-2N} \iint f(x) (|\xi|^2 + \varepsilon^2)^{-p} d^{2N}\xi d^{2N}x. \quad \square \end{aligned}$$

We continue with a technical lemma, in the spirit of [93]. Consider the approximate unit $\{e_K\}_{K \in \mathbb{N}} \subset \mathcal{A}_c$ where $e_K := \sum_{0 \leq |n| \leq K} f_{nn}$. These e_K are projectors with a natural ordering: $e_K \star_\theta e_L = e_L \star_\theta e_K = e_K$ for $K \leq L$, and they are local units for \mathcal{A}_c .

Lemma 9.30. *Let $f \in \mathcal{A}_{c,K}$. Then*

$$\pi^\theta(f) (\mathcal{D}^2 + \varepsilon^2)^{-N} - \pi^\theta(f) \left(\pi^\theta(e_K) (\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e_K) \right)^N \in \mathcal{L}^1(\mathcal{H}).$$

Proof. For simplicity we use the notation $e := e_K$ and $e_n := e_{K+n}$. By the boundedness of $\pi^\theta(f)$, we may assume that $f = e \in \mathcal{A}_{c,K}$.

Because $e_n \star_\theta e = e \star_\theta e_n = e$, it is clear that

$$\pi^\theta(e) (\mathcal{D} + \lambda)^{-1} (1 - \pi^\theta(e_n)) = \pi^\theta(e) (\mathcal{D} + \lambda)^{-1} [\mathcal{D}, \pi^\theta(e_n)] (\mathcal{D} + \lambda)^{-1}. \quad (134)$$

Also, $\pi^\theta(e) [\mathcal{D}, \pi^\theta(e_n)] = [\mathcal{D}, \pi^\theta(e \star_\theta e_n)] - [\mathcal{D}, \pi^\theta(e)] \pi^\theta(e_n) = 0$ because we have the relation $[\mathcal{D}, \pi^\theta(e)] \pi^\theta(e_n) = [\mathcal{D}, \pi^\theta(e)]$ for $n \geq 1$. We obtain

$$\begin{aligned} A_n &:= \pi^\theta(e) (\mathcal{D} + \lambda)^{-1} [\mathcal{D}, \pi^\theta(e_n)] (\mathcal{D} + \lambda)^{-1} \\ &= \pi^\theta(e) (\mathcal{D} + \lambda)^{-1} [\mathcal{D}, \pi^\theta(e_1)] (\mathcal{D} + \lambda)^{-1} [\mathcal{D}, \pi^\theta(e_n)] (\mathcal{D} + \lambda)^{-1} \\ &= \pi^\theta(e) (\mathcal{D} + \lambda)^{-1} [\mathcal{D}, \pi^\theta(e_1)] \pi^\theta(e_2) (\mathcal{D} + \lambda)^{-1} [\mathcal{D}, \pi^\theta(e_n)] (\mathcal{D} + \lambda)^{-1} = \dots \\ &= \left(\pi^\theta(e) (\mathcal{D} + \lambda)^{-1} \right) \left([\mathcal{D}, \pi^\theta(e_1)] (\mathcal{D} + \lambda)^{-1} \right) \left([\mathcal{D}, \pi^\theta(e_2)] (\mathcal{D} + \lambda)^{-1} \right) \dots \\ &\quad \dots \left([\mathcal{D}, \pi^\theta(e_n)] (\mathcal{D} + \lambda)^{-1} \right). \end{aligned}$$

Taking $n = 2N$ here, A_{2N} appears as a product of $2N + 1$ terms in parentheses, each in $\mathcal{L}^{2N+1}(\mathcal{H})$ by Lemma 9.21. Hence, by Hölder's inequality, A_{2N} is trace-class and therefore $\pi^\theta(e)(\mathcal{D} + \lambda)^{-1}(1 - \pi^\theta(e_{2N})) \in \mathcal{L}^1(\mathcal{H})$. Thus,

$$\begin{aligned} & \pi^\theta(e) (\mathcal{D}^2 + \varepsilon^2)^{-1} (1 - \pi^\theta(e_{4N})) \\ &= \pi^\theta(e) (\mathcal{D} - i\varepsilon)^{-1} (1 - \pi^\theta(e_{2N}) + \pi^\theta(e_{2N})) (\mathcal{D} + i\varepsilon)^{-1} (1 - \pi^\theta(e_{4N})) \\ &= \pi^\theta(e) (\mathcal{D} - i\varepsilon)^{-1} (1 - \pi^\theta(e_{2N})) (\mathcal{D} + i\varepsilon)^{-1} (1 - \pi^\theta(e_{4N})) \\ &\quad + \pi^\theta(e) (\mathcal{D} - i\varepsilon)^{-1} \pi^\theta(e_{2N}) (\mathcal{D} + i\varepsilon)^{-1} (1 - \pi^\theta(e_{4N})) \in \mathcal{L}^1(\mathcal{H}). \end{aligned} \quad (135)$$

This is to say $\pi^\theta(e)(\mathcal{D}^2 + \varepsilon^2)^{-1} \sim \pi^\theta(e)(\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e_{4N})$. Shifting this property, we get

$$\begin{aligned} \pi^\theta(e)(\mathcal{D}^2 + \varepsilon^2)^{-N} &\sim \pi^\theta(e)(\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e_{4N})(\mathcal{D}^2 + \varepsilon^2)^{-N+1} \\ &\sim \pi^\theta(e)(\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e_{4N})(\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e_{8N})(\mathcal{D}^2 + \varepsilon^2)^{-N+2} \sim \dots \\ &\sim \pi^\theta(e)(\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e_{4N})(\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e_{8N}) \dots (\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e_{4N^2}). \end{aligned}$$

By identity (130), the last term on the right equals

$$\begin{aligned} & \pi^\theta(e)(\mathcal{D} + i\varepsilon)^{-1} \pi^\theta(e)(\mathcal{D} - i\varepsilon)^{-1} \pi^\theta(e_{4N})(\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e_{8N}) \dots (\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e_{4N^2}) \\ &\quad + \pi^\theta(e)(\mathcal{D} + i\varepsilon)^{-1} [\mathcal{D}, \pi^\theta(e)] (\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e_{4N})(\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e_{8N}) \dots (\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e_{4N^2}). \end{aligned}$$

The last term is trace-class because it is a product of N terms in $\mathcal{L}^p(\mathcal{H})$ for $p > N$ and one term in $\mathcal{L}^q(\mathcal{H})$ for $q > 2N$, by Lemma 9.21. Removing the second $\pi^\theta(e)$ once again, by the ordering property of the local units e_K yields

$$\begin{aligned} & \pi^\theta(e)(\mathcal{D} + i\varepsilon)^{-1} \pi^\theta(e)(\mathcal{D} - i\varepsilon)^{-1} \pi^\theta(e_{4N})(\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e_{8N}) \dots (\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e_{4N^2}) \\ &= \pi^\theta(e)(\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e)(\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e_{8N}) \dots (\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e_{4N^2}) \\ &\quad + \pi^\theta(e)(\mathcal{D}^2 + \varepsilon^2)^{-1} [\mathcal{D}, \pi^\theta(e)] (\mathcal{D} - i\varepsilon)^{-1} \pi^\theta(e_{4N})(\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e_{8N}) \dots (\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e_{4N^2}). \end{aligned}$$

The last term is still trace-class, hence

$$\pi^\theta(e)(\mathcal{D}^2 + \varepsilon^2)^{-N} \sim \pi^\theta(e)(\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e)(\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e_{8N}) \dots (\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e_{4N^2}).$$

This algorithm, applied another $(N - 1)$ times, yields the result:

$$\pi^\theta(e)(\mathcal{D}^2 + \varepsilon^2)^{-N} \sim \left(\pi^\theta(e)(\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e) \right)^N. \quad \square$$

We retain the following consequence.

Corollary 9.31. $\text{Tr}^+ \left(\pi^\theta(g) [\pi^\theta(f), (\mathcal{D}^2 + \varepsilon^2)^{-N}] \right) = 0$ for any $g \in \mathcal{S}$ and any projector $f \in \mathcal{A}_c$.

Proof. This follows from Lemma 9.30 applied to $\pi^\theta(f)(\mathcal{D}^2 + \varepsilon^2)^{-N}$ and its adjoint. \square

Now we are finally ready to evaluate the Dixmier traces.

Proposition 9.32. For $f \in \mathcal{S}$, any Dixmier trace Tr^+ of $\pi^\theta(f) (\mathcal{D}^2 + \varepsilon^2)^{-N}$ is independent of ε , and

$$\text{Tr}^+\left(\pi^\theta(f) (\mathcal{D}^2 + \varepsilon^2)^{-N}\right) = \frac{2^N \Omega_{2N}}{2N(2\pi)^{2N}} \int f(x) d^{2N}x = \frac{1}{N!(2\pi)^N} \int f(x) d^{2N}x.$$

Proof. We will first prove it for $f \in \mathcal{A}_c$. Choose e a unit for f , that is, $e \star_\theta f = f \star_\theta e = f$. By Lemmata 9.28 and 9.30, and because $\mathcal{L}^1(\mathcal{H})$ lies inside the kernel of the Dixmier trace, we obtain

$$\text{Tr}^+\left(\pi^\theta(f) (\mathcal{D}^2 + \varepsilon^2)^{-N}\right) = \text{Tr}^+\left(\pi^\theta(f) (\pi^\theta(e) (\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e))^N\right).$$

Lemma 9.30 applied to $f = e$ implies that $(\pi^\theta(e) (\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e))^N$ is a positive operator in $\mathcal{L}^{1+}(\mathcal{H})$, since it is equal to $\pi^\theta(e) (\mathcal{D}^2 + \varepsilon^2)^{-N}$ plus a term in $\mathcal{L}^1(\mathcal{H})$. Thus, [10, Thm. 5.6] yields (since the limit converges, any Dixmier trace will give the same result):

$$\begin{aligned} \text{Tr}^+\left(\pi^\theta(f) (\mathcal{D}^2 + \varepsilon^2)^{-N}\right) &= \lim_{s \downarrow 1} (s-1) \text{Tr} \left[\pi^\theta(f) (\pi^\theta(e) (\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e))^{Ns} \right] \\ &= \lim_{s \downarrow 1} (s-1) \text{Tr} \left(\pi^\theta(f) \pi^\theta(e) (\mathcal{D}^2 + \varepsilon^2)^{-Ns} \pi^\theta(e) + E_{Ns} \right), \end{aligned} \quad (136)$$

where

$$E_{Ns} := \pi^\theta(f) \left(\pi^\theta(e) (\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e) \right)^{Ns} - \pi^\theta(f) \pi^\theta(e) (\mathcal{D}^2 + \varepsilon^2)^{-Ns} \pi^\theta(e).$$

Lemma 9.30 again shows that $E_N \in \mathcal{L}^1(\mathcal{H})$.

Now for $s > 1$, the first term $\pi^\theta(f) \left(\pi^\theta(e) (\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e) \right)^{Ns}$ of E_{Ns} is in $\mathcal{L}^1(\mathcal{H})$. In effect, using Lemma 9.21 and since $\pi^\theta(e) (\mathcal{D}^2 + \varepsilon^2)^{-1} \in \mathcal{L}^p(\mathcal{H})$ for $p > N$, we get that $\pi^\theta(e) (\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e) \in \mathcal{L}^{Ns}(\mathcal{H})$. This operator being positive, one concludes

$$\left(\pi^\theta(e) (\mathcal{D}^2 + \varepsilon^2)^{-1} \pi^\theta(e) \right)^{Ns} \in \mathcal{L}^1(\mathcal{H}).$$

The second term $\pi^\theta(f) \pi^\theta(e) (\mathcal{D}^2 + \varepsilon^2)^{-Ns} \pi^\theta(e)$ lies in $\mathcal{L}^1(\mathcal{H})$ too, because

$$\|\pi^\theta(e) (\mathcal{D}^2 + \varepsilon^2)^{-Ns} \pi^\theta(e)\|_1 = \|(\mathcal{D}^2 + \varepsilon^2)^{-Ns/2} \pi^\theta(e)\|_2^2 = \|\pi^\theta(e) (\mathcal{D}^2 + \varepsilon^2)^{-Ns/2}\|_2^2$$

is finite by Lemma 9.19. So $E_{Ns} \in \mathcal{L}^1(\mathcal{H})$ for $s \geq 1$, and (136) implies

$$\begin{aligned} \text{Tr}^+\left(\pi^\theta(f) (\mathcal{D}^2 + \varepsilon^2)^{-N}\right) &= \lim_{s \downarrow 1} (s-1) \text{Tr} \left(\pi^\theta(f) \pi^\theta(e) (\mathcal{D}^2 + \varepsilon^2)^{-Ns} \pi^\theta(e) \right) \\ &= \lim_{s \downarrow 1} (s-1) \text{Tr} \left(\pi^\theta(f) (\mathcal{D}^2 + \varepsilon^2)^{-Ns} \right). \end{aligned}$$

Applying now Lemma 9.29, we obtain

$$\begin{aligned} \text{Tr}^+\left(\pi^\theta(f) (\mathcal{D}^2 + \varepsilon^2)^{-N}\right) &= \lim_{s \downarrow 1} (s-1) \text{Tr}(1_{2N}) \text{Tr} \left(L_f^\theta (-\Delta + \varepsilon^2)^{-Ns} \right) \\ &= 2^N (2\pi)^{-2N} \lim_{s \downarrow 1} (s-1) \iint f(x) (|\xi|^2 + \varepsilon^2)^{-Ns} d^{2N}\xi d^{2N}x \\ &= \frac{1}{N!(2\pi)^N} \int f(x) d^{2N}x, \end{aligned}$$

where the identity

$$\int (|\xi|^2 + \varepsilon^2)^{-Ns} d^{2N}\xi = \pi^N \frac{\Gamma(N(s-1))}{\Gamma(Ns) \varepsilon^{2N(s-1)}},$$

and $\Gamma(N\alpha) \sim 1/N\alpha$ as $\alpha \downarrow 0$ have been used. The proposition is proved for $f \in \mathcal{A}_c$.

Finally, take f arbitrary in \mathcal{S} , and recall that $\{e_K\}$ is an approximate unit for \mathcal{A}_θ . Since $f = g \star_\theta h$ for some $g, h \in \mathcal{S}$, Corollary 9.31 implies

$$\begin{aligned} & \left| \text{Tr}^+ \left((\pi^\theta(f) - \pi^\theta(e_K \star_\theta f \star_\theta e_K)) (\mathcal{D}^2 + \varepsilon^2)^{-N} \right) \right| \\ &= \left| \text{Tr}^+ \left((\pi^\theta(f) - \pi^\theta(e_K \star_\theta f)) (\mathcal{D}^2 + \varepsilon^2)^{-N} \right) \right| \\ &= \left| \text{Tr}^+ \left((\pi^\theta(g) - \pi^\theta(e_K \star_\theta g)) \pi^\theta(h) (\mathcal{D}^2 + \varepsilon^2)^{-N} \right) \right| \\ &\leq \|\pi^\theta(g) - \pi^\theta(e_K \star_\theta g)\|_\infty \text{Tr}^+ \left| \pi^\theta(h) (\mathcal{D}^2 + \varepsilon^2)^{-N} \right|. \end{aligned}$$

Since $\|\pi^\theta(g) - \pi^\theta(e_K \star_\theta g)\|_\infty \leq (2\pi\theta)^{-N/2} \|g - e_K \star_\theta g\|_2$ tends to zero when K increases, the proof is complete because $e_K \star_\theta f \star_\theta e_K$ lies in \mathcal{A}_c and

$$\int [e_K \star_\theta f \star_\theta e_K](x) d^{2N}x \rightarrow \int f(x) d^{2N}x \quad \text{as } K \uparrow \infty. \quad \square$$

Remark 9.33. *Similar arguments to those of this section (or a simple comparison argument) show that for $f \in \mathcal{S}$,*

$$\text{Tr}^+ \left(\pi^\theta(f) (|\mathcal{D}| + \varepsilon)^{-2N} \right) = \text{Tr}^+ \left(\pi^\theta(f) (\mathcal{D}^2 + \varepsilon^2)^{-N} \right).$$

In conclusion: the analytical and spectral dimension of Moyal planes coincide. And Lemma 9.28, Proposition 9.32 and the previous remark have concluded the proof of Theorem 9.23.

The conclusion is that $(\mathcal{A}, \tilde{\mathcal{A}}, \mathcal{H}, \mathcal{D}, \chi, J)$ defines a non-compact spectral triple; recall that we already know that both \mathcal{A} and its preferred compactification $\tilde{\mathcal{A}}$ are pre- C^* -algebras.

Theorem 9.34. *The Moyal planes $(\mathcal{A}, \tilde{\mathcal{A}}, \mathcal{H}, \mathcal{D}, J, \chi)$ are connected real non-compact spectral triples of spectral dimension $2N$.*

One can compute the Yang–Mills action (71) of this triple:

Theorem 9.35. *Let $\omega = -\omega^* \in \Omega^1 \mathcal{A}_\theta$. Then the Yang–Mills action $YM(V)$ of the universal connection $\delta + \omega$, with $V = \tilde{\pi}^\theta(\omega)$, is equal to*

$$S_{YM}(V) = -\frac{1}{4g^2} \int F^{\mu\nu} \star_\theta F_{\mu\nu}(x) d^{2N}x = -\frac{1}{4g^2} \int F^{\mu\nu}(x) F_{\mu\nu}(x) d^{2N}x,$$

where $F_{\mu\nu} := \frac{1}{2}(\partial_\mu A_\nu - \partial_\nu A_\mu + [A_\mu, A_\nu]_{\star_\theta})$ and A_μ is defined by $V = L^\theta(A_\mu) \otimes \gamma^\mu$.

The spectral action has been computed in [41, 42]. As one can expected, it is the same, up to few universal coefficients, to the one of Theorem 8.13.

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