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# POLYAKOV'S FORMULATION OF $2 d$ BOSONIC STRING THEORY 

COLIN GUILLARMOU, RÉMI RHODES, AND VINCENT VARGAS


#### Abstract

Using probabilistic methods, we first define Liouville quantum field theory on Riemann surfaces of genus $\mathbf{g} \geqslant 2$ and show that it is a conformal field theory. We use the partition function of Liouville quantum field theory to give a mathematical sense to Polyakov's partition function of noncritical bosonic string theory $[\mathrm{Po}]$ (also called $2 d$ bosonic string theory) and to Liouville quantum gravity. More specifically, we show the convergence of Polyakov's partition function over the moduli space of Riemann surfaces in genus $\mathbf{g} \geqslant 2$ in the case of $D \leqslant 1$ boson. This is done by performing a careful analysis of the behaviour of the partition function at the boundary of moduli space. An essential feature of our approach is that it is probabilistic and non perturbative. The interest of our result is twofold. First, to the best of our knowledge, this is the first mathematical result about convergence of string theories. Second, our construction describes conjecturally the scaling limit of higher genus random planar maps weighted by Conformal Field Theories: we make precise conjectures about this statement at the end of the paper.


## 1. Introduction

In physics, string theory or more generally Euclidean $2 d$ Quantum Gravity (LQG) is an attempt to quantize the Einstein-Hilbert functional coupled to matter fields (matter is replaced by the free bosonic string in the case of string theory). The problematic can be briefly summarized as follows.

First of all, a quantum field theory on a surface $M$ can be viewed as a way to define a measure $e^{-S_{g}(\phi)} D \phi$ over an infinite dimensional space $E$ of fields $\phi$ living over $M$ (typically $\phi$ are sections of some bundles over $M$ ), where $D \phi$ is a "uniform measure" and $S_{g}: E \rightarrow \mathbb{R}$ is a functional on $E$ called the action, depending on a background Riemannian metric $g$ on $M$. The total mass of the measure

$$
\begin{equation*}
Z(g):=\int_{E} e^{-S_{g}(\phi)} D \phi \tag{1.1}
\end{equation*}
$$

is called the partition function. Defining the $n$-points correlation functions amounts to taking $n$ points $x_{1}, \ldots, x_{n} \in M$ and weights $\alpha_{1}, \ldots, \alpha_{n} \in \mathbb{R}$ and to defining

$$
Z\left(g ;\left(x_{1}, \alpha_{1}\right), \ldots,\left(x_{n}, \alpha_{n}\right)\right):=\int_{E} e^{\sum_{i=1}^{n} \alpha_{i} \phi\left(x_{i}\right)} e^{-S_{g}(\phi)} D \phi,
$$

at least if the fields $\phi$ are functions on $M$.
A conformal field theory (CFT in short) on a surface is a quantum field theory which possesses certain conformal symmetries. More specifically, the partition function $Z(g)$ of a CFT satisfies the diffeormorphism invariance $Z\left(\psi^{*} g\right)=Z(g)$ for all smooth diffeomorphism
$\psi: M \rightarrow M$ and a so-called conformal anomaly of the following form: for all $\omega \in C^{\infty}(M)$

$$
\begin{equation*}
Z\left(e^{\omega} g\right)=Z(g) \exp \left(\frac{\mathbf{c}}{96 \pi} \int_{M}\left(|d \omega|_{g}^{2}+2 K_{g} \omega\right) \mathrm{dv}_{g}\right) \tag{1.2}
\end{equation*}
$$

where $\mathbf{c} \in \mathbb{R}$ is called the central charge of the theory, $K_{g}$ is the scalar curvature of $g$ and $\mathrm{dv}_{g}$ the volume form. The $n$-points correlation functions should also satisfy similar types of conformal anomalies and diffeormorphism invariance (see (4.2) and (4.4)). Usually, it is difficult to give a mathematical sense to (1.1) because the measure $D \phi$, which is formally the Lebesgue measure on an infinite dimensional space, does not exist mathematically. Hence, CFT's are mostly studied using axiomatic and algebraic methods, or perturbative methods (formal stationary phase type expansions): see for example [FMS, Ga].

Liouville Quantum Field Theory. The first part of our work is to construct Liouville quantum field theory (LQFT in short) on Riemann surface of genus $\mathbf{g} \geqslant 2$ and to show that this is a CFT. We use probabilistic methods to give a mathematical sense to the path integral (1.1), when $S_{g}(\phi)=S_{L}(g, \phi)$ is the classical Liouville action, a natural convex functional coming from the theory of uniformisation of Riemann surfaces that we describe now. Given a two dimensional connected compact Riemannian manifold ( $M, g$ ) without boundary, we define the Liouville functional on $C^{1}$ maps $\varphi: M \rightarrow \mathbb{R}$ by

$$
\begin{equation*}
S_{L}(g, \varphi):=\frac{1}{4 \pi} \int_{M}\left(|d \varphi|_{g}^{2}+Q K_{g} \varphi+4 \pi \mu e^{\gamma \varphi}\right) \mathrm{dv}_{g} \tag{1.3}
\end{equation*}
$$

where $Q, \mu, \gamma>0$ are parameters to be discussed later. If $Q=\frac{2}{\gamma}$, finding the minimizer $u$ of this functional allows one to uniformize $(M, g)$. Indeed, the metric $g^{\prime}=e^{\gamma u} g$ has constant scalar curvature $K_{g^{\prime}}=-2 \pi \mu \gamma^{2}$ and it is the unique such metric in the conformal class of $g$. The quantization of the Liouville action is precisey LQFT: one wants to make sense of the following finite measure on some appropriate functional space $\Sigma$ (to be defined later) made up of (generalized) functions $\varphi: M \rightarrow \mathbb{R}$

$$
\begin{equation*}
F \mapsto \Pi_{\gamma, \mu}(g, F):=\int_{\Sigma} F(\varphi) e^{-S_{L}(g, \varphi)} D \varphi \tag{1.4}
\end{equation*}
$$

where $D \varphi$ stands for the "formal uniform measure" on $\Sigma$. Up to renormalizing this measure by its total mass, this formalism describes the law of some random (generalized) function $\varphi$ on $\Sigma$, which stands for the ( $\log$-)conformal factor of a random metric of the form $e^{\gamma \varphi} \mathrm{g}$ on $M$. In physics, LQFT is known to be CFT with central charge $\mathbf{c}_{\mathrm{L}}:=1+6 Q^{2}$ continuously ranging in $[25, \infty)$ for the particular values

$$
\begin{equation*}
\gamma \in] 0,2], \quad Q=\frac{2}{\gamma}+\frac{\gamma}{2} . \tag{1.5}
\end{equation*}
$$

Of course, this description is purely formal and giving a mathematical description of this picture is a longstanding problem, which goes back to the work of Polyakov [Po]. The rigorous construction of such an object has been carried out recently in [DKRV] in genus 0, [DRV] in genus 1 (see also [HRV] for the case of the disk). Let us also mention that Duplantier-Miller-Sheffield have developed in [DMS] a theory that lies at the "boundary" of LQFT: with the help of the Gaussian Free Field, their approach gives sense in the case of the sphere, disk or plane to some form of 2-point correlation function of LQFT and the corresponding
equivalence class of random measures (defined up to a non trivial subgroup of biholomorphic transformations of the domain). Yet, another approach by Takhtajan-Teo [TaTe] was to develop a perturbative approach (a semiclassical Liouville theory in the so-called background field formalism approach): in this non-probabilistic approach, LQFT is expanded as a formal power series in $\gamma$ around the minimum of the action (1.3) and the parameter $Q$ in the action is given by its value in classical Liouville theory $Q=\frac{2}{\gamma}$.

We consider the genus $\mathbf{g} \geqslant 2$ case and give a mathematical (and non perturbative) definition to (1.4). To explain the result, we need to summarize the construction. On a compact surface $M$ with genus $\mathbf{g} \geqslant 2$, we fix a smooth metric $g$ on $M$ and define for $s \in \mathbb{R}$ the Sobolev space $H^{s}(M):=\left(1+\Delta_{g}\right)^{-s / 2}\left(L^{2}(M)\right)$ of order $s$ with scalar product defined using the metric $g$. Using the theory of Gaussian free field (GFF in short), we show that for each $s>0$ there is a measure $\mathcal{P}^{\prime}$ on $H^{-s}(M)$ which is independent of the choice of metric $g$ in the conformal class $[g]$, and which represents the following Gaussian formal measure defined for each $F \in L^{1}\left(H^{-s}(M), \mathcal{P}^{\prime}\right)$ by

$$
\begin{equation*}
\int F(\varphi) e^{-\frac{1}{4 \pi} \int_{M}|\nabla \varphi|_{g}^{2} \mathrm{dv}_{g}} D \varphi:=\frac{1}{\sqrt{\operatorname{det}^{\prime}\left(\Delta_{g}\right)}} \int F(\varphi) d \mathcal{P}^{\prime}(\varphi) \tag{1.6}
\end{equation*}
$$

where $\operatorname{det}^{\prime}\left(\Delta_{g}\right)$ is the regularized determinant of the Laplacian, defined as in [RaSi]. The method to do this is to consider a probability space $(\Omega, \mathcal{F}, \mathbb{P})$ and a sequence $\left(a_{j}\right)_{j}$ of i.i.d. real Gaussians $\mathcal{N}(0,1)$ and to consider the random variable (called GFF)

$$
\begin{equation*}
X_{g}=\sqrt{2 \pi} \sum_{j \geqslant 1} a_{j} \frac{\varphi_{j}}{\sqrt{\lambda_{j}}} \tag{1.7}
\end{equation*}
$$

with values in $H^{-s}(M)$, if $\left(\varphi_{j}\right)_{j} \geqslant 0$ is an orthonormal basis of eigenfunctions of $\Delta_{g}$ with eigenvalues $\left(\lambda_{j}\right)_{j \geqslant 0}$ (and with $\lambda_{0}=0$ ). The covariance of $X_{g}$ is the Green function of $\frac{1}{2 \pi} \Delta_{g}$ and there is a probability mesure $\mathcal{P}$ on $H_{0}^{-s}(M):=\left\{u \in H^{-s}(M) ;\langle u, 1\rangle=0\right\}$ so that the law of $X_{g}$ is given by $\mathcal{P}$ and for each $\phi \in H_{0}^{s}(M),\left\langle X_{g}, \phi\right\rangle$ is a random variable on $\Omega$ with zero mean and variance $2 \pi\left\langle\Delta_{g}^{-1} \phi, \phi\right\rangle$. Then $H^{-s}(M)=H_{0}^{-s}(M) \oplus \mathbb{R}$ and we define $\mathcal{P}^{\prime}$ as the pushforward of the measure $\mathcal{P} \otimes d c$ under the mapping $(X, c) \in H_{0}^{-s}(M) \times \mathbb{R} \mapsto c+X$, where $d c$ is the uniform Lebesgue measure in $\mathbb{R}$. The formal equality (1.6) is an analogy with the finite dimensional setting. The next tool needed to the construction is Gaussian multiplicative chaos theory introduced by Kahane [Ka], which allows us to define the random measure $\mathcal{G}_{g}^{\gamma}:=e^{\gamma X_{g}} \mathrm{dv}_{g}$ on $M$ for $0<\gamma \leqslant 2$ when $X_{g}$ is the GFF. This is done by using a renormalisation procedure. We can then define the quantity which plays the role of the formal integral (1.4) as follows: for $F: H^{-s}(M) \rightarrow \mathbb{R}($ with $s>0)$ a bounded continuous functional, we set

$$
\begin{align*}
\Pi_{\gamma, \mu}(g, F):= & \left(\operatorname{det}^{\prime}\left(\Delta_{g}\right) / \operatorname{Vol}_{g}(M)\right)^{-1 / 2}  \tag{1.8}\\
& \times \int_{\mathbb{R}} \mathbb{E}\left[F\left(c+X_{g}\right) \exp \left(-\frac{Q}{4 \pi} \int_{M} K_{g}\left(c+X_{g}\right) \mathrm{dv}_{g}-\mu e^{\gamma c} \mathcal{G}_{g}^{\gamma}(M)\right)\right] d c
\end{align*}
$$

and call it the functional integral of LQFT (when $F=1$ this is the partition function). Our first result is that this quantity is finite and satisfies diffeomorphism invariance and a certain conformal anomaly when $Q=\frac{\gamma}{2}+\frac{2}{\gamma}$.

Theorem 1.1 (LQFT is a CFT). Let $Q=\frac{\gamma}{2}+\frac{2}{\gamma}$ with $\gamma \leqslant 2$ and $g$ be a smooth metric on $M$. For each bounded continuous functional $F: H^{-s}(M) \rightarrow \mathbb{R}$ (with $s>0$ ) and each $\omega \in C^{\infty}(M)$,
$\Pi_{\gamma, \mu}\left(e^{\omega} g, F\right)$ is finite and satisfies the following conformal anomaly:

$$
\Pi_{\gamma, \mu}\left(e^{\omega} g, F\right)=\Pi_{\gamma, \mu}\left(g, F\left(\cdot-\frac{Q}{2} \omega\right)\right) \exp \left(\frac{1+6 Q^{2}}{96 \pi} \int_{M}\left(|d \omega|_{g}^{2}+2 K_{g} \omega\right) \mathrm{dv}_{g}\right)
$$

Let $g$ be any metric on $M$ and $\psi: M \rightarrow M$ be an orientation preserving diffeomorphism, then we have for each bounded measurable $F: H^{-s}(M) \rightarrow \mathbb{R}$ with $s>0$

$$
\Pi_{\gamma, \mu}\left(\psi^{*} g, F\right)=\Pi_{\gamma, \mu}(g, F(\cdot \circ \psi))
$$

As a quantum field theory, the other objects of importance for LQFT are the correlation functions. In Section 4.3, we define the $n$-points correlation functions with vertex operators $e^{\alpha_{i} X_{g}\left(x_{i}\right)}$ where $\alpha_{i}$ are weights and $x_{i} \in M$ some points, and we show their conformal anomaly required to be a CFT. This amounts somehow to taking $F(\varphi)=\prod_{i=1}^{n} e^{\alpha_{i} \varphi\left(x_{i}\right)}$ in (1.8), but it again requires renormalisation since $\varphi$ lives in $H^{-s}(M)$ with $s>0$. At this level, the construction follows the method initiated by [DKRV] on the sphere. We notice that for the sphere, only the $n$-points correlation functions for $n \geqslant 3$ are well defined, while here the partition function is already well-defined.

Liouville Quantum Gravity and Polyakov partition function. Our next result is the main part of the paper and consists in giving a sense to the Liouville quantum gravity (LQG in short) partition function following the work of Polyakov [Po]. We stress that, even though the objects comes from theoretical physics, our result and proof is purely mathematical and can be viewed, from the perspective of a mathematician, as a a way to understand the behaviour of some natural interesting function near the boundary of moduli space, namely the LQFT partition function.

Given a connected closed surface $M$ with genus $\mathbf{g} \geqslant 2$, quantizing the coupling of the gravitational field with matter fields amounts to making sense of the formal integral (partition function),

$$
\begin{equation*}
Z=\int_{\mathcal{R}} e^{-S_{\mathrm{EH}}(g)}\left(\int e^{-S_{\mathrm{M}}\left(g, \phi_{m}\right)} D_{g} \phi_{m}\right) D g \tag{1.9}
\end{equation*}
$$

where the measure $D g$ runs over the space of Riemannian structures $\mathcal{R}$ on $M$ (the space of metrics $g$ modulo diffeomorphisms), the functional integral for matter fields $\int e^{-S_{\mathrm{M}}\left(g, \phi_{m}\right)} D_{g} \phi_{m}$ stands for the quantization of an action $\phi_{m} \mapsto S_{\mathrm{M}}\left(g, \phi_{m}\right)$ over an infinite dimensional space of fields describing matter, and $S_{\text {EH }}$ is the Einstein-Hilbert action

$$
\begin{equation*}
S_{\mathrm{EH}}(g)=\frac{1}{2 \kappa} \int_{M} K_{g} d \mathrm{v}_{g}+\mu_{0} \operatorname{Vol}_{g}(M) \tag{1.10}
\end{equation*}
$$

where $\kappa$ is the Einstein constant, $\mu_{0} \in \mathbb{R}$ is the cosmological constant. The measure $D g$ represents the formal Riemannian measure associated to the $L^{2}$ metric on the space of Riemannian metrics (its reduction to $\mathcal{R}$ ). There are several possible choices for the matter fields and we shall focus on the choice described in Polyakov [Po]. Giving a mathematical definition to the functional integral (1.9) has been a real challenge, and Polyakov [Po] suggested a decomposition of this integral in the case of bosonic string theory with $D$ free bosons. In that case, $Z_{M}(g):=\int e^{-S_{\mathrm{M}}\left(g, \phi_{m}\right)} D_{g} \phi_{m}$ is the partition function of a CFT with central charge $\mathbf{c}_{\mathrm{M}}=D$, and is mathematically given by a certain power of the determinant of the Laplacian. The argument of Polyakov [Po] (see also D'Hoker-Phong [DhPh]) for defining (1.9) was based on the
observation that each metric $g$ on $M$ can be decomposed as

$$
\begin{equation*}
g=\psi^{*}\left(e^{\omega} g_{\tau}\right) \tag{1.11}
\end{equation*}
$$

where $\omega \in C^{\infty}(M), \psi$ is a diffeomorphism and $\left(g_{\tau}\right)_{\tau \in \mathcal{M}_{\mathrm{g}}}$ is a family of hyperbolic metrics on $M$ parameterizing the moduli space $\mathcal{M}_{g}$ of genus-g surfaces. We recall that $\mathcal{M}_{g}$ is the space of equivalence classes of conformal structures: it is a $6 \mathbf{g}-6$ dimensional orbifold equipped with a natural metric, called the Weil-Petersson metric, whose volume form denoted $d \tau$ has finite volume. In this way, the space of Riemannian structures $\mathcal{R}$ is identified to the product of moduli space $\mathcal{M}_{\mathrm{g}}$ with the Weyl group $C^{\infty}(M)$ acting on metrics by $(\varphi, g) \mapsto e^{\varphi} g$. Applying the change of variables (1.11) in the formal integral (1.9) produces a Jacobian (called Ghost determinant) taking into account the quotient of the space of metrics by the space of diffeomorphisms of $M$. The Ghost determinant turns out to be the partition function of a CFT with central charge $\mathbf{c}_{\text {ghost }}=-26$ and Polyakov noticed that at the specific value $D=26$ the conformal anomaly of $Z_{\mathrm{M}}$ cancels out that of the Ghost term, giving rise to a Weyl invariant partition function

$$
\begin{equation*}
Z=\int_{\mathcal{M}_{g}} Z_{M}\left(g_{\tau}\right) Z_{\text {Ghost }}\left(g_{\tau}\right) \sqrt{\operatorname{det} J_{g_{\tau}}} d \tau . \tag{1.12}
\end{equation*}
$$

called critical string theory ${ }^{1}$. Concretely, this was further discussed by D'Hoker-Phong [DhPh] who wrote

$$
\begin{equation*}
Z_{\text {Ghost }}(g)=\left(\frac{\operatorname{det}\left(P_{g}^{*} P_{g}\right)}{\operatorname{det} J_{g}}\right)^{1 / 2}, \quad Z_{\mathrm{M}}(g)=C\left(\frac{\operatorname{det}^{\prime}\left(\Delta_{g}\right)}{\operatorname{Vol}_{g}(M)}\right)^{-\frac{D}{2}} \tag{1.13}
\end{equation*}
$$

for some constant $C$, where the determinants are defined by zeta regularizations, $P_{g}$ is a firstorder elliptic operator mapping 1-forms to trace-free symmetric 2-tensors, and $J_{g}$ is the Gram matrix of a fixed basis of ker $P_{g}^{*}$ (see Section 5.1 further details). Then Belavin-Knizhnik [BeKn] and Wolpert [Wo2] proved that the integral (1.12) with $D=26$ diverges at the boundary of (the compactification of) moduli space, a problematic fact in order to establish well-posedness of the partition function for critical $D=26$ (bosonic) strings.

Noncritical string theories are not formulated within the critical dimension $D=26$, yet they are Weyl invariant. The idea, emerging once again from the paper [Po], is that for $D \neq 26$ the integral (1.9) possesses one further degree of freedom to be integrated over corresponding to the Weyl factor $e^{\omega}$ in (1.11). For $D \leqslant 1$, hence $\mathbf{c}_{\mathrm{M}} \leqslant 1$, Polyakov argued that integrating this factor requires using Liouville quantum field theory. In other words, applying once again the change of variables (1.11) to (1.9) yields

$$
\begin{equation*}
Z=\int_{\mathcal{M}_{\mathbf{g}}} Z_{M}\left(g_{\tau}\right) Z_{\text {Ghost }}\left(g_{\tau}\right) \Pi_{\gamma, \mu}\left(g_{\tau}, 1\right) \sqrt{\operatorname{det} J_{g_{\tau}}} d \tau \tag{1.14}
\end{equation*}
$$

where $\Pi_{\gamma, \mu}\left(g_{\tau}, 1\right)$ is the partition function of LQFT in the background metric $g_{\tau}$. As explained above, the partition function $\Pi_{\gamma, \mu}\left(g_{\tau}, 1\right)$ depends on two parameters $\gamma$ and $\mu$ (some of which are related in some specific way to $\mathbf{c}_{\mathrm{M}}$ ). Later, Polyakov's argument was generalized by David and Distler-Kawai [Da, DiKa] to CFT type matter field theories with central charge $\mathbf{c}_{\mathrm{M}} \leqslant 1$ (thus including the case $D=1$ ). In that CFT context, the integral (1.14) is often called partition function of Liouville quantum gravity. This approach has an important consequence related to

[^0]string theory as it paves the way to a rigorous construction of noncritical bosonic string theory ${ }^{2}$ provided one can make sense of (1.14). This is the main purpose of this paper. The importance of this theory is discussed in great details in [Polch, section 5.1] or [Kleb] for instance.

Once having defined the partition function $\Pi_{\gamma, \mu}(g, 1)$ of LQFT, we introduce the partition function for LQG

$$
\begin{equation*}
Z_{\mathrm{LQG}}:=\int_{\mathcal{M g}_{\mathrm{g}}} Z_{\mathrm{M}}\left(g_{\tau}\right) Z_{\mathrm{Ghost}}\left(g_{\tau}\right) \Pi_{\gamma, \mu}\left(g_{\tau}, 1\right) \sqrt{\operatorname{det} J_{g_{\tau}}} d \tau \tag{1.15}
\end{equation*}
$$

where we choose for the matter partition function

$$
\begin{equation*}
Z_{\mathrm{M}}(g)=\left(\frac{\operatorname{det}^{\prime} \Delta_{g}}{\operatorname{Vol}_{g}(M)}\right)^{-\frac{\mathrm{c}_{\mathrm{M}}}{2}}, \tag{1.16}
\end{equation*}
$$

the ghost determinant $Z_{\text {Ghost }}\left(g_{\tau}\right)$ being defined by (1.13) and $d \tau$ the Weil-Petersson measure. Notice that (1.16) is nothing but the partition function (1.13) for $D=c_{\mathrm{M}}$ free bosons extended to all possible values $D \leqslant 1$. The parameters in (1.8) are tuned in such a way that the global conformal anomaly of the product

$$
Z_{\mathrm{M}}\left(g_{\tau}\right) Z_{\text {Ghost }}\left(g_{\tau}\right) \Pi_{\gamma, \mu}\left(g_{\tau}, 1\right)
$$

vanishes, hence ensuring Weyl invariance of the whole theory (1.15). In view of Theorem 1.1, this gives the relation

$$
\mathbf{c}_{\mathrm{M}}-26+1+6 Q^{2}=0,
$$

hence determining the value of $\gamma($ encoded by $Q$ ) in terms of the central charge of the matter fields. We refer to Section 5.1 for more explanations.

The main result of this paper is the following:
Theorem 1.2 (Convergence of the partition function). For surfaces of genus $\mathbf{g} \geqslant 2$, the integral defining the partition function $Z_{\mathrm{LQG}}$ of (1.15) converges for $\left.\left.\gamma \in\right] 0,2\right]$, that is for $\mathbf{c}_{\mathrm{M}} \leqslant 1$.

The integral defining $Z_{\mathrm{LQG}}$ in the case $\mathbf{c}_{\mathrm{M}}=0$ corresponds to the case of pure gravity (i.e. no matter), $\mathbf{c}_{\mathrm{M}}=-2$ to uniform spanning trees and $\mathbf{c}_{\mathrm{M}}=1$ (equivalently $D=1$ ) to noncritical strings (or $D=2$ string theory). As far as we know, this is the first proof of convergence of string theory on hyperbolic surfaces with fixed genus.

Using this Theorem, we can now see the metric $g$ on $M$ as a random variable with law ruled by the partition function (1.15). The Riemannian volume and modulus of this random metric are called the quantum gravity volume form (LQG volume form) and quantum gravity modulus, see Theorem 5.1. Furthermore we formulate conjectures relating the LQG volume form to the scaling limit of random planar maps in the case of pure gravity $\mathbf{c}_{\mathrm{M}}=0$, hence providing the scaling limit of the model studied in [Mi], or to the scaling limit of random planar maps weighted by the discrete Gaussian free field in the case $\mathbf{c}_{\mathrm{M}}=1$, see Section 5.5. Though we do not explicitly write a conjecture, we further mention here that the limit of random planar maps with fixed topology weighted by uniform spanning trees corresponds to $\mathbf{c}_{\mathrm{M}}=-2$.

[^1]The main input of our paper is the proof of Theorem 1.2. We need to analyse the integrands in (1.15) near the boundary of moduli space and show that we can control them. The moduli space $\mathcal{M}_{\mathrm{g}}$ can be viewed as a $6 \mathbf{g}-6$ dimensional non-compact orbifold of hyperbolic metrics on $M$, that can be compactified in such a way that its boundary corresponds to pinching closed geodesics. Hyperbolic metrics on $M$ corresponding to points in $\partial \mathcal{M}_{\mathrm{g}}$ are complete hyperbolic surfaces with cusps and finite volume. The estimates of Wolpert [Wo2] describe the parts involving the ghost and matter terms. The heart of our work is to analyse the behaviour of Gaussian multiplicative chaos under degeneracies of the hyperbolic surfaces: this is rather involved since there are in general small eigenvalues of Laplacian tending to 0 and the covariance of the GFF (i.e. the Green function) is thus diverging. There is yet a huge conceptual gap between the cases $\mathbf{c}_{\mathrm{M}}<1$ and $\mathbf{c}_{\mathrm{M}}=1$. Roughly, the reason is the following: the product $Z_{\mathrm{M}}\left(g_{\tau}\right) Z_{\text {Ghost }}\left(g_{\tau}\right)$ is at leading order comparable to $\prod_{j} e^{-\frac{\pi^{2}}{3 \ell_{j}}\left(1-\frac{\mathbf{c}_{\mathrm{M}}}{2}\right)}$, where the product runs over pinched geodesics with lengths $\ell_{j} \rightarrow 0$ when approaching the boundary of the (compactification of) moduli space - see subsection 2.3 for more precise statementswhereas $\Pi_{\gamma, \mu}\left(g_{\tau}, 1\right)$ is comparable to $\prod_{j} e^{-\frac{\pi^{2}}{3_{j}}\left(-\frac{1}{2}\right)} \times F\left(\mathcal{G}_{g}^{\gamma}(M)\right)$ where $F\left(\mathcal{G}_{g}^{\gamma}(M)\right)$ is an explicit functional expectation of the Gaussian multiplicative chaos $\mathcal{G}_{g}^{\gamma}(M)$. Hence, for $\mathbf{c}_{\mathrm{M}}<1$, we prove soft estimates on the functional $F\left(\mathcal{G}_{g}^{\gamma}(M)\right)$ that are enough to get an exponential decay of the product $Z_{\mathrm{M}}\left(g_{\tau}\right) Z_{\text {Ghost }}\left(g_{\tau}\right) \Pi_{\gamma, \mu}\left(g_{\tau}, 1\right)$ at the boundary of $\mathcal{M}_{\mathrm{g}}$ and thus integrability with respect to the Weil-Petersson measure. In the case $\mathbf{c}_{\mathrm{M}}=1$, the leading exponential behaviors cancel out exactly so that the analysis must determine the polynomial corrections behind the leading exponential behavior, rendering the computations more much intricate. In order to analyse the mass of the Gaussian multiplicative chaos measure in these degenerating regions, we need to prove very sharp estimates (that, as far as we know, are new) on the Green function and on the eigenfunctions associated to the small eigenvalues in the pinched necks of the surface, as functions of the moduli space parameters $\tau$ when $\tau$ approaches the boundary $\partial \mathcal{M}_{\mathrm{g}}$. Roughly speaking, the crucial observation is that the GFF behaves like two independent Brownian motions in the variable transverse to the closed geodesic being pinched, and this allows us to translate the problem in terms of explicit functionals of Brownian motion. For (and only for) $\mathbf{c}_{\mathrm{M}}=1$, we show that the pinching produces an extra rate of decay of $\Pi_{\gamma, \mu}\left(g_{\tau}, 1\right)$ as we approach $\partial \mathcal{M}_{\mathbf{g}}$, implying the convergence.

To conclude this introduction, we point out that an interesting different approach to define path integrals for random Kähler metrics on surfaces was introduced recently by Ferrari-Klevtsov-Zelditch [FKZ, KlZe], but the link with our work is not established rigorously.

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## 2. Geometric background and Green functions

2.1. Uniformisation of compact surfaces of genus $\mathrm{g} \geqslant 2$. Let $M$ be a compact surface of genus $\mathbf{g} \geqslant 2$ and let $g$ be a smooth Riemannian metric. Recall that Gauss-Bonnet tells us
that

$$
\begin{equation*}
\int_{M} K_{g} \mathrm{dv}_{g}=4 \pi \chi(M) \tag{2.1}
\end{equation*}
$$

where $\chi(M)=(2-2 \mathbf{g})$ is the Euler characteristic, $K_{g}$ the scalar curvature of $g$ and $\mathrm{dv}_{g}$ the Riemannian measure. The uniformisation theorem says that in the conformal class

$$
[g]:=\left\{e^{\varphi} g ; \varphi \in C^{\infty}(M)\right\}
$$

of $g$, there exists a unique metric $g_{0}=e^{\varphi_{0}} g$ of scalar curvature $K_{g_{0}}=-2$. For a metric $\hat{g}=e^{\varphi} g$, one has the relation

$$
K_{\hat{g}}=e^{-\varphi}\left(\Delta_{g} \varphi+K_{g}\right)
$$

where $\Delta_{g}=d^{*} d$ is the non-negative Laplacian (here $d$ is exterior derivative and $d^{*}$ its adjoint). Finding $\varphi_{0}$ is achieved by minimizing the following functional

$$
F: C^{\infty}(M) \rightarrow \mathbb{R}^{+}, \quad F(\varphi):=\int_{M}\left(\frac{1}{2}|d \varphi|_{g}^{2}+K_{g} \varphi+2 e^{\varphi}\right) \mathrm{dv}_{g}
$$

and taking $\varphi_{0}$ to be the function where $F(\varphi)$ is minimum at $\varphi=\varphi_{0}$. We will embed this functional into a more general one, depending on three parameters, called Liouville functional: let $\gamma, Q, \mu>0$ and define

$$
\begin{equation*}
S_{L}(g, \varphi):=\frac{1}{4 \pi} \int_{M}\left(|d \varphi|_{g}^{2}+Q K_{g} \varphi+4 \pi \mu e^{\gamma \varphi}\right) \mathrm{dv}_{g} . \tag{2.2}
\end{equation*}
$$

When $Q=2 / \gamma$ and $\pi \mu \gamma^{2}=1$, we can write $S_{L}(g, \varphi)=\frac{1}{2 \gamma^{2} \pi} F(\gamma \varphi)$. In fact, if $\hat{g}=e^{\omega} g$ for some $\varphi$, the functional $S_{L}$ satisfies the relation

$$
S_{L}\left(\hat{g}, \varphi-\frac{\omega}{\gamma}\right)=S_{L}(g, \varphi)+\frac{1}{4 \pi} \int_{M}\left(\left(\frac{1}{\gamma^{2}}-\frac{Q}{\gamma}\right)|d \omega|_{g}^{2}-\frac{Q}{\gamma} K_{g} \omega+\left(Q-\frac{2}{\gamma}\right) \varphi \Delta_{g} \omega\right) \mathrm{dv}_{g}
$$

and in particular if $Q=2 / \gamma$ it satisfies

$$
\begin{equation*}
S_{L}\left(\hat{g}, \varphi-\frac{\omega}{\gamma}\right)=S_{L}(g, \varphi)-\frac{1}{4 \pi \gamma^{2}} \int_{M}\left(|d \omega|_{g}^{2}+2 K_{g} \omega\right) \mathrm{dv}_{g}, \tag{2.3}
\end{equation*}
$$

which is sometimes called conformal anomaly: changing the conformal factor of the metric entails a variation of the functional proportional to the Liouville functional. Similar properties will be shared by the quantum version of the Liouville theory, which fall under the scope of Conformal Field Theory (CFT for short, more later). At this stage let us just mention that we will show that the value of $Q$ for the quantum Liouville theory to possess a conformal anomaly has to be adjusted to take into account quantum effects. More precisely we will have in the quantum theory

$$
\begin{equation*}
Q=\frac{2}{\gamma}+\frac{\gamma}{2} . \tag{2.4}
\end{equation*}
$$

2.2. Hyperbolic surfaces, Teichmüller space and Moduli space. Let $M$ be a surface of genus $\mathbf{g} \geqslant 2$. The set of smooth metrics on $M$ is a Fréchet manifold denoted by $\operatorname{Met}(M)$ and contained in the Fréchet space of smooth symmetric tensors $C^{\infty}\left(M, S^{2} M\right)$ of order 2. This space has a natural $L^{2}$ metric given by

$$
\begin{equation*}
\left\langle h_{1}, h_{2}\right\rangle:=\int_{M}\left\langle h_{1}, h_{2}\right\rangle_{g} \operatorname{dv}_{g} \tag{2.5}
\end{equation*}
$$

where $h_{1}, h_{2} \in T_{g} \operatorname{Met}(M)=C^{\infty}\left(M, S^{2} M\right)$ and $\langle\cdot, \cdot\rangle_{g}$ is the usual scalar product on endomorphisms of $T M$ when we identify symmetric 2 -tensors with endomorphisms of $T M$ through the metric $g$. A metric with Gaussian curvature -1 will be called hyperbolic, we denote by $\operatorname{Met}_{-1}(M)$ the set of such metrics on $M$. The group $\mathcal{D}(M)$ of smooth diffeomorphisms acts smoothly and properly on $\operatorname{Met}(M)$ and on $\operatorname{Met}_{-1}(M)$ by pull-back $\phi . g:=\phi^{*} g$, moreover it acts by isometries with respect to the metric (2.5). The subgroup $\mathcal{D}_{0}(M) \subset \mathcal{D}(M)$ of elements contained in the connected component of the Identity also acts properly and smoothly and $\operatorname{Mod}(M):=\mathcal{D}(M) / \mathcal{D}_{0}(M)$ is a discrete subgroup called mapping class group or moduli group. The Fréchet space $C^{\infty}(M)$ acts on $\operatorname{Met}(M)$ by conformal multiplication $(\varphi, g) \mapsto e^{\varphi} g$. The orbits of this action are called conformal classes and the conformal class of a metric $g$ is denoted by $[g]$.

The Teichmüller space of $M$ is defined by

$$
\mathcal{T}(M):=\operatorname{Met}_{-1}(M) / \mathcal{D}_{0}(M)
$$

By taking slices transverse to the action of $\mathcal{D}_{0}(M)$, we can put a structure of smooth manifold with real dimension $6 \mathbf{g}-6$, it is topologically a ball, and its tangent space at a metric $g$ (representing a class in $\mathcal{T}(M)$ ) can be identified naturally with the space of divergence-free trace-free tensors with respect to $g$ by choosing appropriately the slice. Teichmüller space has a complex structure and is equipped with a natural Kähler metric called the Weil-Petersson metric, which is defined by

$$
\left\langle h_{1}, h_{2}\right\rangle_{\mathrm{WP}}:=\int_{M}\left\langle h_{1}^{\mathrm{tf}}, h_{2}^{\mathrm{tf}}\right\rangle_{g} \mathrm{dv}_{g}
$$

if $h_{1}, h_{2} \in T_{g} \mathcal{T}(M)$ and $h_{i}^{\mathrm{tf}}=h_{i}-\frac{1}{2} \operatorname{Tr}_{g}\left(h_{i}\right) g$ denotes the trace-free part. The Weil-Petersson metric is not the metric induced by (2.5) after quotienting by $\mathcal{D}_{0}(M)$ but it is rather induced by the $L^{2}$ metric on almost complex structures, when we identify almost complex structures with metrics of constant curvature. We refer for to the book of Tromba [Tr] for more details about this approach of Teichmüller theory, the Weil-Petersson metric is discussed in Section 2.6 there.

The group $\operatorname{Mod}(M)$ acts properly discontinuously on $\mathcal{T}(M)$ by isometries of the WeilPetersson metric, but the action is not free and there are elements of finite order. The quotient $\mathcal{M}(M):=\mathcal{T}(M) / \operatorname{Mod}(M)$ is a Riemannian orbifold called the moduli space of $M$, its orbifold singularities corresponding to hyperbolic metrics admitting isometries. Since $\mathcal{T}(M), \operatorname{Mod}(M)$ and $\mathcal{M}(M)$ actually depend only on the genus $\mathbf{g}$ of $M$, we shall denote them $\mathcal{T}_{\mathbf{g}}, \operatorname{Mod}_{\mathbf{g}}$ and $\mathcal{M}_{\mathrm{g}}$. The manifold $\mathcal{M}_{\mathrm{g}}$ is open but can be compactified into $\overline{\mathcal{M}_{\mathrm{g}}}$, the locus of the compactification is a divisor $D \subset \overline{\mathcal{M}_{\mathrm{g}}}$ and the Weil-Petersson distance is complete on that space. Since we will need to understand the behavior of certain quantities on the moduli space, we now recall its geometry near the divisor $D$ and we shall follow the description given by Wolpert ([Wo1, Wo2, Wo3]) for this compactification. On a surface $M$ of genus $\mathbf{g}$, there is a unique geodesic in each free homotopy class, and we call a partition of $M$ a collection of $3 \mathbf{g}-3$ simple closed curves $\left\{\gamma_{1}, \ldots, \gamma_{3 g-3}\right\}$ which are not null-homotopic and not mutually homotopic. If $g \in \operatorname{Met}_{-1}(M)$, there is a unique simple geodesic homotopic to each $\gamma_{j}$ and we obtain a decomposition of $(M, g)$ into $2 \mathrm{~g}-2$ hyperbolic pants (a pant is a topological sphere with 3 disk removed, equipped with a hyperbolic metric and with totally geodesic boundary). A subpartition of $M$ is a collection of $n_{p}$ simple curves $\left\{\gamma_{1}, \ldots, \gamma_{n_{p}}\right\}$ which are not null-homotopic
and not mutually homotopic, with $n_{p} \leqslant 3 \mathbf{g}-3$; they disconnect the surface into surfaces with boundary. A surface $\left(M_{0}, g_{0}\right)$ in $\partial \overline{\mathcal{M}_{\mathbf{g}}}$ is a surface with nodes: $M_{0}$ is the interior of a compact surface $M$ with $n_{p}$ simple curves $\gamma_{1}, \ldots, \gamma_{n_{p}}$ removed and $g_{0}$ is a complete hyperbolic metric with finite volume on $M_{0}$, the metric in a collar neighborhood $[-1,1]_{\rho} \times(\mathbb{R} / \mathbb{Z})_{\theta}$ of each $\gamma_{j}$ (with $\gamma_{j}=\{\rho=0\}$ ) being

$$
g_{0}=\frac{d \rho^{2}}{\rho^{2}}+\rho^{2} d \theta^{2}
$$

Notice that these corresponds to a pair of hyperbolic cusps, each one isometric to ( $\mathbb{R}_{t}^{+} \times$ $\left.(\mathbb{R} / \mathbb{Z})_{\theta}, d t^{2}+e^{-2 t} d \theta^{2}\right)$ by setting $x= \pm e^{-t}$. Now there is a neighborhood $\bar{U}_{g_{0}}$ of $g_{0}$ in $\overline{\mathcal{M}_{\mathbf{g}}}$ represented by hyperbolic metrics $g_{s, \tau}$ on $M_{0}$ for some parameter $(s, \tau) \in \mathbb{C}^{3 g-3-m} \times \mathbb{C}^{m}$ near 0 , with $g_{s, \tau}$ which are smooth metrics on $M$ when $\tau \neq 0$ and complete smooth metrics with hyperbolic cusps on $M_{0}$ when $\tau=0$, and $g_{0,0}=g_{0}$. Moreover, the metrics $g_{s, \tau}$ are continuous with respect to $(s, \tau)$ on compact sets of $M_{0}$ for $(s, \tau)$ near 0 (in the $C^{\infty}$ topology), they are given in the fixed collar neighborhood $[-1,1]_{\rho} \times(\mathbb{R} / \mathbb{Z})_{\theta}$ of $\gamma_{j}$ by

$$
\begin{equation*}
g_{s, \tau}=e^{\varphi_{s, \tau}}\left(\frac{d \rho^{2}}{\varepsilon_{j}^{2}+\rho^{2}}+\left(\varepsilon_{j}^{2}+\rho^{2}\right) d \theta^{2}\right) \tag{2.6}
\end{equation*}
$$

with $\varepsilon_{j}:=2 \pi^{2} /|\log | \tau_{j}| |$ and $\varphi_{s, \tau} \in C^{\infty}(M)$ satisfies

$$
e^{\varphi_{s, \tau}}-1 \rightarrow 0 \text { as }(s, \tau) \rightarrow 0
$$

in $C^{0}$ norm (and in fact in $C^{\infty}$ on compact sets of $M_{0}$ ). Here we notice that the metric $\frac{d \rho^{2}}{\varepsilon_{j}^{2}+\rho^{2}}+\left(\varepsilon_{j}^{2}+\rho^{2}\right) d \theta^{2}$ has curvature -1 in the collar and is isometric to

$$
\begin{equation*}
\left[-t_{j}, t_{j}\right]_{t} \times(\mathbb{R} / \mathbb{Z})_{\theta}, \quad d t^{2}+\varepsilon_{j}^{2} \cosh (t)^{2} d \theta^{2} \tag{2.7}
\end{equation*}
$$

by setting $\varepsilon_{j} \sinh (t)=\rho$ and $\varepsilon_{j} \sinh \left(t_{j}\right)=1$. The geodesic $\gamma_{j}\left(g_{s, t}\right)$ for $g_{s, t}$ homotopic to $\gamma_{j}$ has length

$$
\ell_{j}\left(g_{s, \tau}\right)=\varepsilon_{j}(1+o(1)) \text { as }|(s, \tau)| \rightarrow 0
$$

and is contained in a neighborhood $\left[-c \varepsilon_{j}, c \varepsilon_{j}\right] \times \mathbb{R} / \mathbb{Z}$ of the collar near $\gamma_{j}$ for some $c>0$ independent of $\varepsilon_{j}$ (or equivalently in $t \in[-c, c]$ ). Using $\left|e^{\varphi_{s, \tau}}-1\right|<\delta$ for some small $\delta>0$, the set $B_{j}=\left\{m \in M ; d_{g}\left(\gamma_{j}\left(g_{s, \tau}\right), m\right) \geqslant\left|\log \ell_{j}\left(g_{s, \tau}\right)\right|\right\}$ is contained in a compact set of $M_{0}$ uniform in $(s, \tau)$ where the metrics depend continously on $(s, \tau)$. We can then use geodesic normal coordinates with respect to $g$ around $\gamma_{j}\left(g_{s, \tau}\right)$ and the collar $\mathcal{C}_{j}\left(g_{s, \tau}\right)=M \backslash B_{j}$ is isometric to

$$
\begin{equation*}
\left[-\log \ell_{j}\left(g_{s, \tau}\right), \log \ell_{j}\left(g_{s, \tau}\right)\right]_{t} \times(\mathbb{R} / \mathbb{Z})_{\theta}, \quad d t^{2}+\ell_{j}\left(g_{s, \tau}\right)^{2} \cosh (t)^{2} d \theta^{2} \tag{2.8}
\end{equation*}
$$

To summarize, the geometry is uniformly bounded outside $\cup_{j=1}^{r} \mathcal{C}_{j}(g)$ for metrics $g$ in a small neighborhood $\bar{U}_{g_{0}}$ of $g_{0}$ in $\overline{\mathcal{M}_{\mathbf{g}}}$.

The loops $\left(\gamma_{j}\right)_{j}$ define a subpartition. The open strata of $D$ correspond to subpartitions up to equivalence by elements in $\operatorname{Mod}_{\mathbf{g}}$. For each $\beta>0$, the set of metrics in $\mathcal{M}_{\mathbf{g}}$ such that all geodesics have length larger than $\beta$ is a compact subset of $\mathcal{M}_{\mathbf{g}}$ called the $\beta$-thick part. The $\beta$-thin part of $\mathcal{M}_{\mathbf{g}}$ is the complement of the $\beta$-thick part. By Lemma 6.1 of [Wo2], there exists a constant $\beta>0$ so that the $\beta$-thin part of $\mathcal{M}_{\mathbf{g}}$ is covered by a finite set of neighborhoods $U\left(\mathrm{SP}_{j}\right), j=1, \ldots, n_{p}$ where $\mathrm{SP}_{j}$ denote some subpartitions of $M$ and $U\left(\mathrm{SP}_{j}\right)$ denote the set of surfaces in Teichmüller space (up to $\operatorname{Mod}_{g}$ equivalence) for which the geodesics in the homotopy
class of curves of $\mathrm{SP}_{j}$ have length less than $\beta$ and the other ones have length bounded below by $\beta / 2$. Each $U\left(\mathrm{SP}_{j}\right)$ is a neighborhood of a strata of $D$.

For each pants decomposition of the surface (with genus $\mathbf{g}$ ), one has associated coordinates $\tau=\left(\ell_{1}, \ldots, \ell_{3 \mathrm{~g}-3}, \theta_{1}, \ldots, \theta_{3 \mathrm{~g}-3}\right)$ where $\ell_{j}$ are the lengths of the simple closed geodesics bounding the pair of pants and $\theta_{j} \in[0,2 \pi)$ are the twist angles (see [Wo1]). The Weil-Petersson volume form is given in these coordinates by

$$
\begin{equation*}
d \tau:=C_{\mathbf{g}} \prod_{j=1}^{3 \mathrm{~g}-3} \ell_{j} d \theta_{j} d \ell_{j} \tag{2.9}
\end{equation*}
$$

for some constant $C_{\mathrm{g}}>0$ depending only on the genus.
2.3. Determinant of Laplacians. For a Riemannian metric $g$ on a connected oriented compact surface $M$, the non-negative Laplacian $\Delta_{g}=d^{*} d$ has discrete spectrum $\operatorname{Sp}\left(\Delta_{g}\right)=\left(\lambda_{j}\right)_{j \in \mathbb{N}_{0}}$ with $\lambda_{0}=0$ and $\lambda_{j} \rightarrow+\infty$. We can define the determinant of $\Delta_{g}$ by

$$
\operatorname{det}^{\prime}\left(\Delta_{g}\right)=\exp \left(-\left.\partial_{s} \zeta(s)\right|_{s=0}\right)
$$

where $\zeta(s):=\sum_{j=1}^{\infty} \lambda_{j}^{-s}$ is the spectral zeta function of $\Delta_{g}$, which admits a meromorphic continuation from $\operatorname{Re}(s) \gg 1$ to $s \in \mathbb{C}$ and is holomorphic at $s=0$. We recall that if $\hat{g}=e^{\varphi} g$ for some $\varphi \in C^{\infty}(M)$, one has the so-called Polyakov formula (see [OPS, eq. (1.13)])

$$
\begin{equation*}
\log \frac{\operatorname{det}^{\prime}\left(\Delta_{\hat{g}}\right)}{\operatorname{Vol}_{\hat{g}}(M)}=\log \frac{\operatorname{det}^{\prime}\left(\Delta_{g}\right)}{\operatorname{Vol}_{g}(M)}-\frac{1}{48 \pi} \int_{M}\left(|d \varphi|_{g}^{2}+2 K_{g} \varphi\right) \operatorname{dv}_{g} \tag{2.10}
\end{equation*}
$$

where $K_{g}$ is the scalar curvature of $g$ as above. It is interesting to compare with the conformal anomaly (2.3) of the Liouville action $S_{L}$. To compute $\operatorname{det}^{\prime}\left(\Delta_{g}\right)$, it thus suffices to know it for an element in the conformal class, and by the uniformisation theorem we can choose a metric $g$ of scalar curvature -2 (or equivalently Gaussian curvature -1 ) if $M$ has genus $\mathbf{g} \geqslant 2$. Such hyperbolic surface can be realized as a quotient $\Gamma \backslash \mathbb{H}^{2}$ of the hyperbolic half-plane

$$
\mathbb{H}^{2}:=\{z \in \mathbb{C} ; \operatorname{Im}(z)>0\} \text { with metric } g_{\mathbb{H}^{2}}=\frac{|d z|^{2}}{(\operatorname{Im}(z))^{2}}
$$

by a discrete co-compact subgroup $\Gamma \subset \mathrm{PSL}_{2}(\mathbb{R})$ with no torsion. In each free homotopy class on $M=\Gamma \backslash \mathbb{H}^{2}$, there is a unique closed geodesic, and we can form the Selberg zeta function

$$
Z_{g}(s)=\prod_{\gamma \in \mathcal{P}} \prod_{k=0}^{\infty}\left(1-e^{-(s+k) \ell(\gamma)}\right), \quad \operatorname{Re}(s)>1
$$

where $\mathcal{P}$ denotes the set of primitive closed geodesics of $(M, g) \simeq \Gamma \backslash \mathbb{H}^{2}$ and $\ell(\gamma)$ are their lengths (recall that primitive closed geodesics are oriented closed geodesics that are not iterates of another closed geodesic). By the work of Selberg, The function $Z_{g}(s)$ admits an analytic continuation to $s \in \mathbb{C}$ and it is proved by D'Hoker-Phong [DhPh] that

$$
\begin{equation*}
\operatorname{det}^{\prime} \Delta_{g}=Z_{g}^{\prime}(1) e^{(2 \mathbf{g}-2) C} \tag{2.11}
\end{equation*}
$$

where $C$ is an explicit universal constant (see also Sarnak [Sa]). The behavior of $Z_{g}^{\prime}(1)$ near the boundary of $\mathcal{M}_{\mathbf{g}}$ is studied by Wolpert [Wo2]: there exists $C_{\mathbf{g}}>0$ a constant depending
only on the genus such that for all $g \in \mathcal{M}_{\mathbf{g}}$

$$
\begin{equation*}
C_{\mathbf{g}}^{-1} \prod_{j=1}^{n_{p}} \frac{e^{-\frac{\pi^{2}}{3 \ell_{j}(g)}}}{\ell_{j}(g)} \prod_{\lambda_{k}(g)<1 / 4} \lambda_{k}(g) \leqslant Z_{g}^{\prime}(1) \leqslant C_{\mathbf{g}} \prod_{j=1}^{n_{p}} \frac{e^{-\frac{\pi^{2}}{3 \ell_{j}(g)}}}{\ell_{j}(g)} \prod_{\lambda_{k}(g)<1 / 4} \lambda_{k}(g) \tag{2.12}
\end{equation*}
$$

where $\lambda_{k}(g)$ are the eigenvalues of $\Delta_{g}$ and $\ell_{j}(g)$ are the lengths of closed geodesics with length less than $\varepsilon>0$ for some small fixed $\varepsilon>0$.

There is another operator which appears in the work of Polyakov [Po] and whose determinant is important in 2D quantum gravity. Let $P_{g}$ be the differential operator mapping differential 1 -forms on $M$ to symmetric trace-free 2 -tensors, defined by

$$
P_{g} \omega:=2 \mathcal{S} \nabla^{g} \omega-\operatorname{Tr}_{g}\left(\mathcal{S} \nabla^{g} \omega\right) g .
$$

If $\nabla^{g}$ is Levi-Civita connection for $g, \operatorname{Tr}_{g}$ denotes the trace with respect to $g$ and $\mathcal{S}$ denotes the orthogonal projection on symmetric 2-tensors. The kernel of $P_{g}$ is the space of conformal Killing vector fields, which is thus trivial in genus $\mathbf{g} \geqslant 2$. Its adjoint $P_{g}^{*}$ is given by $P_{g}^{*} u:=$ $\delta_{g}(u)=-\operatorname{Tr}_{g}\left(\nabla^{g} u\right)$ and called the divergence operator on symmetric trace-free 2 -tensors. Its kernel has real dimension $6 \mathbf{g}-6$ and is conformally invariant. We denote by ( $\phi_{1}, \ldots, \phi_{6 \mathbf{g}-6}$ ) a fixed basis of ker $P_{g}^{*}$ and by $J_{g}$ the matrix $\left(J_{g}\right)_{i j}=\left\langle\phi_{i}, \phi_{j}\right\rangle_{g}$ The operator $P_{g}^{*} P_{g}$ is an elliptic positive self-adjoint second order differential operator acting on 1-forms, and we can define its determinant by

$$
\operatorname{det}\left(P_{g}^{*} P_{g}\right)=\exp \left(-\left.\partial_{s} \zeta_{1}(s)\right|_{s=0}\right), \quad \zeta_{1}(s)=\sum_{j=1}^{\infty} \mu_{j}^{-s}
$$

where $\mu_{j}>0$ are the non-zero eigenvalues of $P_{g}^{*} P_{g}$. The conformal anomaly for this operator is proved by Alvarez [Al, Eq. 4.27] and reads

$$
\begin{equation*}
\log \frac{\operatorname{det}\left(P_{\hat{g}}^{*} P_{\hat{g}}\right)}{\operatorname{det} J_{\hat{g}}}=\log \frac{\operatorname{det}\left(P_{g}^{*} P_{g}\right)}{\operatorname{det} J_{g}}-\frac{13}{24 \pi} \int_{M}\left(|d \varphi|_{g}^{2}+2 K_{g} \varphi\right) \operatorname{dv}_{g} \tag{2.13}
\end{equation*}
$$

if $\hat{g}=e^{\varphi} g$. By [DhPh], one has for $(M, g)$ a hyperbolic surface realized as $\Gamma \backslash \mathbb{H}^{2}$

$$
\begin{equation*}
\operatorname{det}\left(P_{g}^{*} P_{g}\right)^{\frac{1}{2}}=Z_{g}(2) e^{(2 \mathbf{g}-2) C^{\prime}} \tag{2.14}
\end{equation*}
$$

for some universal constant $C^{\prime}$, and $Z_{g}(s)$ is again the Selberg zeta function. The behavior of $Z_{g}(2)$ near the boundary of $\mathcal{M}_{\mathbf{g}}$ is also studied by Wolpert [Wo2]: there exists $C_{\mathbf{g}}>0$ a constant depending only on the genus such that for all $g \in \mathcal{M}_{\mathbf{g}}$

$$
\begin{equation*}
C_{\mathbf{g}}^{-1} \prod_{j=1}^{n_{p}} \frac{e^{-\frac{\pi^{2}}{3 \ell_{j}(g)}}}{\ell_{j}^{3}(g)} \leqslant Z_{g}(2) \leqslant C_{\mathbf{g}} \prod_{j=1}^{n_{p}} \frac{e^{-\frac{\pi^{2}}{3 \ell_{j}(g)}}}{\ell_{j}^{3}(g)} \tag{2.15}
\end{equation*}
$$

where the $\ell_{j}(g)$ are the lengths of closed geodesics with length less than $\varepsilon>0$ for some small fixed $\varepsilon>0$.
2.4. Green function and resolvent of Laplacian. Each compact Riemannian surface $(M, g)$ has a (non-negative) Laplace operator $\Delta_{g}=d^{*} d$ and a Green function $G_{g}$ defined to be the integral kernel of the resolvent operator $R_{g}: L^{2}(M) \rightarrow L^{2}(M)$ satisfying $\Delta_{g} R_{g}=\operatorname{Id}-\Pi_{0}$ where $\Pi_{0}$ is the orthogonal projection in $L^{2}\left(M, \mathrm{dv}_{g}\right)$ on ker $\Delta_{g}$ (the constants). By integral kernel, we mean that for each $f \in L^{2}(M)$

$$
R_{g} f(x)=\int_{M} G_{g}\left(x, x^{\prime}\right) f\left(x^{\prime}\right) \operatorname{dv}_{g}\left(x^{\prime}\right)
$$

It is well-known (see for example [Pa]) that the hyperbolic space $\mathbb{H}^{2}$ also has a family of Green functions (here $d_{\mathbb{H}^{2}}\left(z, z^{\prime}\right)$ denotes the hyperbolic distance between $\left.z, z^{\prime}\right)$

$$
\begin{equation*}
G_{\mathbb{H}^{2}}\left(\lambda ; z, z^{\prime}\right)=F_{\lambda}\left(d_{\mathbb{H}^{2}}\left(z, z^{\prime}\right)\right), \quad \lambda \in D(0,1 / 4) \subset \mathbb{C} \tag{2.16}
\end{equation*}
$$

so that $F_{\lambda}(r)$ is a holomorphic function of $\lambda$ for $r \in(0, \infty)$ satisfying

$$
\begin{equation*}
F_{\lambda}(r) \sim-\frac{1}{2 \pi} \log (r) \text { as } r \rightarrow 0, \quad F_{0}(r)=-\frac{1}{2 \pi} \log (r)+m\left(r^{2}\right) \tag{2.17}
\end{equation*}
$$

with $m$ being a smooth functions on $[0, \infty)$, and $G_{\mathbb{H}^{2}}(s)$ satisfies

$$
\left(\Delta_{\mathbb{H}^{2}}-\lambda\right) G_{\mathbb{H}^{2}}\left(\lambda ; \cdot, z^{\prime}\right)=\delta_{z^{\prime}}
$$

where $\delta_{z^{\prime}}$ denotes the Dirac mass at $z^{\prime}$; in other words, $G_{\mathbb{H}^{2}}(\lambda)$ is the Schwartz kernel of the operator $\left(\Delta_{\mathbb{H}^{2}}-\lambda\right)^{-1}$ on $L^{2}\left(\mathbb{H}^{2}\right)$. To obtain the Green function $G_{g}\left(x, x^{\prime}\right)$, it suffices to know it for $g$ hyperbolic (ie. $g$ has constant Gaussian curvature -1 ) since for any other conformal metric $\hat{g}=e^{\varphi} g$, we have that

$$
\begin{gather*}
G_{\hat{g}}\left(x, x^{\prime}\right)=G_{g}\left(x, x^{\prime}\right)+\alpha-u(x)-u\left(x^{\prime}\right), \\
\text { with } \alpha=\frac{1}{\operatorname{Vol}_{\hat{g}}(M)^{2}}\left\langle G_{g}, 1 \otimes 1\right\rangle_{\hat{g}}, \quad u(x):=\frac{1}{\operatorname{Vol}_{\hat{g}}(M)} \int_{M} G_{g}(x, y) \operatorname{dv}_{\hat{g}}(y) . \tag{2.18}
\end{gather*}
$$

This follows from an easy computation and the identity $\Delta_{\hat{g}}=e^{-\varphi} \Delta_{g}$.
Let us then assume that $g$ is hyperbolic. We have
Lemma 2.1. If $g$ is a hyperbolic metric on the surface $M$, the Green function $G_{g}\left(x, x^{\prime}\right)$ for $\Delta_{g}$ has the following form near the diagonal

$$
\begin{equation*}
G_{g}\left(x, x^{\prime}\right)=-\frac{1}{2 \pi} \log \left(d_{g}\left(x, x^{\prime}\right)\right)+m_{g}\left(x, x^{\prime}\right) \tag{2.19}
\end{equation*}
$$

for some smooth function $m_{g}$ on $M \times M$. Near each point $x_{0} \in M$, there are isothermal coordinates $z$ so that $g=|d z|^{2} / \operatorname{Im}(z)^{2}$ and near $x_{0}$

$$
G_{g}\left(z, z^{\prime}\right)=-\frac{1}{2 \pi} \log \left|z-z^{\prime}\right|+F\left(z, z^{\prime}\right)
$$

with $F$ smooth. Finally, if $\hat{g}$ is any metric conformal to $g$, (2.19) holds with $\hat{g}$ replacing $g$ but with $m_{\hat{g}}$ continuous.

Proof. Near each point $x_{0} \in M$, there is an isometry from a geodesic ball $B_{g}\left(x_{0}, \varepsilon\right)$ for $g$ to the hyperbolic ball $B_{\mathbb{H}^{2}}(0, \varepsilon)$ in $\mathbb{H}^{2}$ ( 0 denotes the center of $\mathbb{H}^{2}$ viewed as the unit disk), which provides in particular some local complex coordinates $z$ near $x_{0}$ so that $g=e^{\varphi}|d z|^{2}$ in the ball $B_{g}\left(x_{0}, \varepsilon\right)$ for some function $\varphi$. In these coordinates,

$$
\begin{equation*}
\log d_{g}\left(x, x^{\prime}\right)=\log d_{\mathbb{H}^{2}}\left(z, z^{\prime}\right)=\log \left|z-z^{\prime}\right|+L\left(z, z^{\prime}\right) \text { with } L \text { smooth } \tag{2.20}
\end{equation*}
$$

where $d_{g}$ denote the distance for the metric $g$. Near any given point $x^{\prime} \in M$, one has

$$
\begin{equation*}
\left(\Delta_{g}-\lambda\right) F_{\lambda}\left(d_{g}\left(\cdot, x^{\prime}\right)\right)-\delta_{x^{\prime}} \in C^{\infty}\left(B_{g}\left(x^{\prime}, \varepsilon\right)\right) \tag{2.21}
\end{equation*}
$$

where $B_{g}\left(x^{\prime}, \varepsilon\right)$ is a geodesic ball of center $x^{\prime}$ and radius $\varepsilon>0$ small. Denote by $R_{g}(\lambda)=$ $\left(\Delta_{g}-\lambda\right)^{-1}$ the resolvent of $\Delta_{g}$ for $\lambda \notin \operatorname{Sp}\left(\Delta_{g}\right)$. By the spectral theorem, at $\lambda=\lambda_{0}$ with $\lambda_{0} \in \operatorname{Sp}\left(\Delta_{g}\right)$ we have the Laurent expansion

$$
R_{g}(\lambda)=\frac{\Pi_{\lambda_{0}}}{\lambda-\lambda_{0}}+R_{g}\left(\lambda_{0}\right)+\mathcal{O}\left(\left(\lambda-\lambda_{0}\right)\right), \quad \lambda \rightarrow \lambda_{0}
$$

for some bounded operator $R_{g}\left(\lambda_{0}\right)$ and $\Pi_{\lambda_{0}}$ being the orthogonal projector on $\operatorname{ker}\left(\Delta_{g}-\lambda_{0}\right)$. Thus we obtain

$$
\left(\Delta_{g}-\lambda_{0}\right) R_{g}\left(\lambda_{0}\right)=\operatorname{Id}-\Pi_{\lambda_{0}}
$$

and by elliptic regularity and (2.21), the Schwartz kernel $G_{g}\left(\lambda ; x, x^{\prime}\right)$ of $R_{g}(\lambda)$ for $\lambda \notin \operatorname{Sp}\left(\Delta_{g}\right)$ satisfies for $d_{g}\left(x, x^{\prime}\right)<\varepsilon$ with $\varepsilon>0$ small enough

$$
\begin{equation*}
G_{g}\left(\lambda ; x, x^{\prime}\right)=F_{\lambda}\left(d_{g}\left(x, x^{\prime}\right)\right)+E_{g}\left(\lambda ; x, x^{\prime}\right) \tag{2.22}
\end{equation*}
$$

with $E_{g}$ some smooth function on $M \times M$ depending meromorphically of $\lambda$. At $\lambda=0$ we deduce (2.19). The part about $\hat{g}$ just follows from (2.18) and the fact that $d_{\hat{g}}\left(x, x^{\prime}\right)=e^{\varphi(x) / 2} d_{g}\left(x, x^{\prime}\right)+$ $\mathcal{O}\left(d_{g}\left(x, x^{\prime}\right)^{2}\right)$ as $x^{\prime} \rightarrow x$.

The function $x \mapsto m_{g}(x, x)$ is often called the Robin constant at $x$. Notice that if we view the hyperbolic metric $g$ as an element representing a point of $\mathcal{T}_{\mathbf{g}}$ and if $\psi \in \operatorname{Mod}_{\mathbf{g}}$, then we have the modular invariance

$$
\begin{equation*}
G_{\psi^{*} g}\left(\lambda ; x, x^{\prime}\right)=G_{g}\left(\lambda ; \psi(x), \psi\left(x^{\prime}\right)\right), \quad m_{\psi^{*} g}\left(x, x^{\prime}\right)=m_{g}\left(\psi(x), \psi\left(x^{\prime}\right)\right) . \tag{2.23}
\end{equation*}
$$

We shall need to describe the Green function $G_{g}$ when the metric $g$ approaches the boundary of the compactification of $\mathcal{M}_{\mathrm{g}}$. It turns out that positive small eigenvalues of $\Delta_{g}$ appear sometime when $g$ approaches a point in $\partial \overline{\mathcal{M}}_{\mathbf{g}}$ : Schoen-Wolpert-Yau [SWY] proved that there exist two positive constants $\alpha_{1}, \alpha_{2}$ depending only on the genus $\mathbf{g}$ so that the $n$-th positive eigenvalue $\lambda_{n}(g)$ of $\Delta_{g}$ satisfy

$$
\alpha_{1} L_{n}(M, g) \leqslant \lambda_{n}(g) \leqslant \alpha_{2} L_{n}(M, g) \text { if } n \leqslant 2 \mathbf{g}-2, \quad \text { and } \alpha_{1} \leqslant \lambda_{2 \mathbf{g}-1} \leqslant \alpha_{2}
$$

where $L_{n}(M, g)$ is the minimum (over subpartitions) of the sums of lengths of simple geodesics in subpartitions of $M$ disconnecting $M$ into $n+1$ connected components. Each metric $g_{0} \in \partial \overline{\mathcal{M}_{\mathbf{g}}}$ is in a stratum corresponding to a subpartition SP containing $n_{p}$ curves, with $n_{s} \leqslant n_{p}$ of these simple curves $\gamma_{1}, \ldots, \gamma_{n_{s}}$ in the subpartition that disconnect the surface $M$ into $m+1$ connected components. There is $c_{0}>0$ depending on $g_{0}$ such that for all $\delta>0$ small enough, there is a neighborhood $\bar{U}_{g_{0}} \subset \overline{\mathcal{M}_{\mathrm{g}}}$ of $g_{0}$ such that for all $g$ in the interior $U_{g_{0}}:=\bar{U}_{g_{0}} \cap \mathcal{M}_{\mathrm{g}}$, there is at most $m$ positive eigenvalues less than $\delta$ and all other eigenvalues are bigger than $c_{0}$. We call these eigenvalues the small eigenvalues of $g$ near $g_{0}$.

Proposition 2.2. Let $\left(M_{0}, g_{0}\right) \in \partial \overline{\mathcal{M}_{\mathbf{g}}}$ where $M_{0}$ is a surface with nodes. For $\delta>0$ arbitrarily small, take $g$ in a sufficiently small open neighborhood $\bar{U}_{g_{0}}$ of $g_{0}$ in $\overline{\mathcal{M}_{\mathbf{g}}}$ so that the small
eigenvalues of $g$ in $U_{g_{0}}=\bar{U}_{g_{0}} \cap \mathcal{M}_{\mathbf{g}}$ satisfy $\lambda_{1}(g) \leqslant \ldots \leqslant \lambda_{m}(g) \leqslant \delta$. The Green function $G_{g}$ restricted to $M_{0}$ can be written for $g \in U_{g_{0}}$ as

$$
\begin{equation*}
G_{g}\left(x, x^{\prime}\right)=\sum_{\lambda_{j}(g) \leqslant \delta} \frac{\Pi_{\lambda_{j}(g)}\left(x, x^{\prime}\right)}{\lambda_{j}(g)}+A_{g}\left(x, x^{\prime}\right) \tag{2.24}
\end{equation*}
$$

where $\Pi_{\lambda_{j}(g)}$ is the orthogonal projector on the corresponding eigenspace. In each compact set $\Omega$ of $M_{0}$, the map $\left(g, x, x^{\prime}\right) \mapsto A_{g}\left(x, x^{\prime}\right)$ is continuous on $\bar{U}_{g_{0}} \times\left(\Omega_{\text {diag }}^{2}\right)$ if $\Omega_{\text {diag }}^{2}:=(\Omega \times \Omega) \backslash$ diag and, near the diagonal of $\Omega \times \Omega$, one has

$$
A_{g}\left(x, x^{\prime}\right)=-\frac{1}{2 \pi} \log \left(d_{g}\left(x, x^{\prime}\right)\right)+B_{g}\left(x, x^{\prime}\right)
$$

with $\left(g, x, x^{\prime}\right) \mapsto B_{g}\left(x, x^{\prime}\right) \in C^{0}\left(\bar{U}_{g_{0}} \times \Omega \times \Omega\right)$. The Schwartz kernel $\sum_{j=1}^{m} \Pi_{\lambda_{j}(g)}\left(x, x^{\prime}\right)$ extends continuously to $\left(g, x, x^{\prime}\right) \in \bar{U}_{g_{0}} \times \Omega \times \Omega$ with value at $g=g^{\prime} \in \partial \bar{U}_{g_{0}}$ the orthogonal projector $\Pi_{0}\left(g^{\prime} ; x, x^{\prime}\right)$ onto $\operatorname{ker}_{L^{2}} \Delta_{g^{\prime}}$.

Proof. After possibly splitting $\Omega$ in smaller pieces, we can assume that the radius of injectivity of all $g \in U_{g_{0}}$ on $\Omega$ is bounded below by some uniform $\alpha>0$. Using the residue formula applied to $R_{g}(\lambda) / \lambda$ in a disk $D(0, \delta)$ of radius $\delta$ centered at $\lambda=0$, one has

$$
\begin{equation*}
R_{g}(0)-\sum_{\lambda_{j}(g) \leqslant \delta} \frac{\Pi_{\lambda_{j}(g)}}{\lambda_{j}(g)}=\frac{1}{2 \pi i} \int_{\partial D(0, \delta)} \frac{R_{g}(\lambda)}{\lambda} d \lambda . \tag{2.25}
\end{equation*}
$$

and we denote by $A_{g}\left(x, x^{\prime}\right)$ the Schwartz kernel of $\frac{1}{2 \pi i} \int_{\partial D(0, \delta)} \frac{R_{g}(\lambda)}{\lambda} d \lambda$. Let $\Omega^{\prime} \subset M_{0}$ be a small neighborhood $\Omega^{\prime}$ of $\Omega$ so that the radius of injectivity of each $g \in U_{g_{0}}$ is bounded below by $\alpha / 2$. Let $L_{g}(\lambda)$ be the operator on $\Omega^{\prime}$ with Schwartz kernel

$$
F_{\lambda}\left(d_{g}\left(x, x^{\prime}\right)\right)
$$

where $F_{\lambda}$ is the function of (2.16). Take $\chi, \widetilde{\chi} \in C_{c}^{\infty}\left(M_{0}\right)$ equal to 1 on $\Omega$ but with support contained in $\Omega^{\prime}$, and such that $\tilde{\chi} \chi=\chi$. Then on $\Omega^{\prime}$ and on $M_{0}$, we have

$$
\begin{equation*}
\left(\Delta_{g}-\lambda\right) \widetilde{\chi} L_{g}(\lambda) \chi=\chi+\left[\Delta_{g}, \widetilde{\chi} L_{g}(\lambda) \chi\right. \tag{2.26}
\end{equation*}
$$

Then multiplying (2.26) by $\chi R_{g}(\lambda)$ on the left, we get

$$
\chi R_{g}(\lambda) \chi=\chi L_{g}(\lambda) \chi-\chi R_{g}(\lambda)\left[\Delta_{g}, \widetilde{\chi}\right] L_{g}(\lambda) \chi .
$$

The operators $\left[\Delta_{g}, \widetilde{\chi}\right\rceil L_{g}(\lambda) \chi$ have smooth kernel (we use that $\left[\Delta_{g}, \widetilde{\chi}\right]=0$ on $\operatorname{supp}(\chi)$ ), and extends continuously to $g \in \bar{U}_{g_{0}}$ since $g$ extends continuously as a smooth metric to $\Omega$ and $d_{g}$ on $\Omega \times \Omega$ as well. Now we use the fact that for $\lambda \in \partial D(0, \delta), g \mapsto R_{g}(\lambda)$ extends continuously to $\bar{U}_{g_{0}}$ as bounded operators $H_{\text {comp }}^{k}\left(M_{0}\right) \rightarrow H_{\text {loc }}^{k}\left(M_{0}\right)$ for all $k \geqslant 0$ by a result of Schulze [Sc]: this implies that $\chi R_{g}(\lambda)\left[\Delta_{g}, \widetilde{\chi}\right] L_{g}(\lambda) \chi$ extend continuously in $g \in \bar{U}_{g_{0}}$ as a family of bounded operators $H^{-k}\left(M_{0}\right) \rightarrow H_{\text {comp }}^{k}\left(M_{0}\right)$ for all $k \geqslant 0$, since $\left[\Delta_{g}, \widetilde{\chi}\right] L_{g}(\lambda) \chi$ maps $H^{-k}\left(M_{0}\right) \rightarrow H^{k}\left(M_{0}\right)$ uniformly in $g \in U_{g_{0}}$. Thus the Schwartz kernels of the operators $\left[\Delta_{g}, \widetilde{\chi}\right] L_{g}(\lambda) \chi$ extend as a uniform family of continuous Schwartz kernels (when $g \in U_{g_{0}}$ ). We then deduce that

$$
\frac{1}{2 \pi i} \int_{\partial D(0, \delta)} \frac{\chi R_{g}(\lambda) \chi}{\lambda} d \lambda=\frac{1}{2 \pi i} \int_{\partial D(0, \delta)} \frac{\chi L_{g}(\lambda) \chi}{\lambda} d \lambda+B_{g}^{\prime}
$$

where $B_{g}^{\prime}$ is a family of operators, with Schwartz kernel $B_{g}^{\prime}\left(x, x^{\prime}\right)$ continuous as a function of $\left(g, x, x^{\prime}\right) \in \bar{U}_{g_{0}} \times \Omega \times \Omega$. Next, since by Cauchy formula $\frac{1}{2 \pi i} \int_{\partial D(0, \delta)} \frac{F_{\lambda}(z)}{\lambda} d \lambda=F_{0}(z)$, we deduce that

$$
\left(\frac{1}{2 \pi i} \int_{\partial D(0, \delta)} \frac{\chi L_{g}(\lambda) \chi}{\lambda} d \lambda\right)\left(x, x^{\prime}\right)=\chi(x) \chi\left(x^{\prime}\right) F_{0}\left(d_{g}\left(x, x^{\prime}\right)\right)
$$

and this Schwartz kernel has the desired property by using (2.17). This ends the proof of (2.24). The proof of the fact that $\sum_{j=1}^{m} \Pi_{\lambda_{j}(g)}\left(x, x^{\prime}\right)$ converge to the projector onto the kernel of $g^{\prime}$ as $g \rightarrow g^{\prime} \in \partial \bar{U}_{g_{0}}$ is essentially the same as what we did (and even simpler) by applying the residue formula to $R_{g}(\lambda)$ in $D(0, \delta)$ instead of $R_{g}(\lambda) / \lambda$. The convergence in $C^{0}$ norm is clear since convergence in $L^{2}$ of $\sum_{j=1}^{m} \Pi_{\lambda_{j}(g)}$ implies convergence in $C^{\infty}$ on $\Omega$ by elliptic regularity.
2.5. Small eigenvalues and associated eigenvectors. In this section, we recall the asymptotics of the small positive eigenvalues $\lambda_{1}(g) \leqslant \ldots \leqslant \lambda_{m}(g)$ as $g \in \mathcal{M}_{\mathrm{g}}$ approaches an element $g_{0} \in \overline{\mathcal{M}_{\mathbf{g}}}$ by following Burger [Bu1, Bu2] and we will see that the proof of [Bu2] also gives an approximation of the projectors $\Pi_{\lambda_{j}(g)}$. Let $\left(M_{0}, g_{0}\right)$ be a surface with nodes, with corresponding subpartition of the closed Riemann surface $M$ given by simple curves $\gamma_{1}, \ldots, \gamma_{n_{p}}$ and $\gamma_{1}, \ldots, \gamma_{n_{s}}$ (with $n_{s} \leqslant n_{p}$ ) are disconnecting $M$ into $m+1$ connected components $S_{1}, \ldots, S_{m+1}$. For all $\delta>0$ small enough, there is a small neighborhood $\bar{U}_{g_{0}} \subset \overline{\mathcal{M}_{\mathbf{g}}}$ of $g_{0}$ in $\overline{\mathcal{M}_{\mathbf{g}}}$ so that for each $g \in U_{g_{0}}=\bar{U}_{g_{0}} \cap \mathcal{M}_{\mathrm{g}}$ there are $m$ small eigenvalues $0<\lambda_{1}(g) \leqslant \ldots \leqslant \lambda_{m}(g) \leqslant \delta$ and all others are larger than a constant $c_{0}>0$ depending only on $g_{0}$. Each metric $g \in U_{g_{0}}$ has a unique simple closed geodesic homotopic to $\gamma_{j}$ for $j \leqslant r$, with length $\ell_{j}(g) \leqslant c_{1} \delta$, while all other primitive closed geodesics have length bigger than $c_{2}>0$, where $c_{1}, c_{2}$ are constants depending only on $g_{0}$. Define the length $L_{i j}(g):=\sum_{k \in E_{i j}} \ell_{k}(g)$ where $E_{i j}=\left\{1 \leqslant k \leqslant n_{s} ; \gamma_{k} \in \partial S_{i} \cap \partial S_{j}\right\}$. Let $\|\cdot\|_{g}$ be the norm on $\mathbb{R}^{m+1}$ given by

$$
\begin{equation*}
\|a\|_{g}^{2}=\sum_{j=1}^{m+1} \operatorname{Vol}_{g}\left(S_{j}\right) a_{j}^{2}, \quad \text { with } a=\left(a_{1}, \ldots, a_{m+1}\right) \tag{2.27}
\end{equation*}
$$

and let $Q_{g}$ be the quadratic form on $\mathbb{R}^{m+1}$ given by

$$
\begin{equation*}
Q_{g}(a)=\sum_{1 \leqslant i, j \leqslant m+1}\left(a_{i}-a_{j}\right)^{2} L_{i j}(g) . \tag{2.28}
\end{equation*}
$$

Notice that $\operatorname{Vol}_{g}\left(S_{j}\right)$ are positive constants depending only on the topology of $S_{j}$ (and not on $g$ ) by Gauss-Bonnet theorem. Then Burger [Bu2] showed the following estimate:

Theorem 2.3 (Burger 1990). If $\nu_{1}(g) \leqslant \ldots \leqslant \nu_{m}(g)$ are the positive eigenvalues of $Q_{g}$ with respect to the norm $\|\cdot\|_{g}$ on $\mathbb{R}^{m+1}$, then there is $C>0$ such that for all $g \in U_{g_{0}}$ and each $1 \leqslant j \leqslant m$

$$
\frac{\nu_{j}(g)}{\pi}\left(1-C \delta^{\frac{1}{2}}\right) \leqslant \lambda_{j}(g) \leqslant \frac{\nu_{j}(g)}{\pi}(1+C \delta|\log \delta|) .
$$

Each simple small geodesic $\gamma_{j}(g)$ of $g$ (homotopic to $\gamma_{j}$ ) has a collar neighborhood

$$
\begin{equation*}
\mathcal{C}_{j}(g)=\left\{x \in M ; \sinh \left(d_{g}\left(x, \gamma_{j}(g)\right)\right) \leqslant 1 / \sinh \left(\ell_{j}(g)\right)\right\} \tag{2.29}
\end{equation*}
$$

and these collars are disjoints one from the other. The set $M \backslash \cup_{j \leqslant n_{s}} \mathcal{C}_{j}(g)$ has $m+1$ connected components $S_{1}^{\prime}, \ldots, S_{m+1}^{\prime}$ that identify naturally with $S_{1}, \ldots, S_{m+1}$. One can define a map

$$
\begin{equation*}
a \in \mathbb{R}^{m+1} \mapsto f_{a} \in H^{1}(M) \tag{2.30}
\end{equation*}
$$

by setting $f_{a}(x)=a_{j}$ if $x \in S_{j}^{\prime}$ and $f_{a}$ being the unique harmonic function in $\mathcal{C}_{j}(g)$ so that $f_{a}$ is continuous on $M$. In [Bu2], Burger proved the following

Lemma 2.4 (Burger 90). There is $C>0$ such that for all $a \in \mathbb{R}^{m+1}$ and all $g \in U_{g_{0}}$

$$
\frac{1}{\pi} Q_{g}(a) \leqslant\left\|d f_{a}\right\|_{L^{2}(M, g)}^{2} \leqslant \frac{1}{\pi} Q_{g}(a)(1+C \delta), \quad\|a\|_{g}^{2}(1-C \delta|\log \delta|) \leqslant\left\|f_{a}\right\|_{L^{2}(M, g)}^{2} \leqslant\|a\|_{g}^{2}
$$

An estimate for $\lambda_{j}(g)$ in terms of the pinched geodesics $\ell_{k}(g)$ is given by Schoen-Wolpert-Yau [SWY]: let $D$ be an $n$-disconnect, i.e a collection of closed simple geodesics $\gamma_{1}(g), \ldots, \gamma_{n_{s}}(g)$ with respective lenghts $\ell_{1}(g), \ldots, \ell_{n_{s}}(g)$ disconnecting $M$ into $n$ connected components, and define $L_{n}(D, g):=\sum_{j=1}^{n_{s}} \ell_{j}(g)$. We set

$$
L_{n}(M, g):=\min _{D \in \mathcal{D}_{n}} L_{n}(D, g)
$$

where $\mathcal{D}_{n}$ is the set of all $n$-disconnects of $M$. Then, ordering the eigenavlues by increasing order, it is proved in [SWY] that there is $C>1$ depending only on the genus of $M$ such that for each $n \leqslant 2 \mathbf{g}-2$

$$
C^{-1} L_{n}(M, g) \leqslant \lambda_{n}(g) \leqslant C L_{n}(M, g)
$$

As a consequence, in a neighborhood $U_{g_{0}} \subset \mathcal{M}_{\mathbf{g}}$ of a metric $g_{0} \in \partial \mathcal{M}_{\mathbf{g}}$, we have the rough estimate for all $j \leqslant m$ and $g \in U_{g_{0}}$

$$
\begin{equation*}
\lambda_{j}(g) \geqslant C^{-1} \ell_{j}(g) \tag{2.31}
\end{equation*}
$$

where $m+1$ is the number of connected components of the surface with cusps $\left(M_{0}, g_{0}\right), n_{s}$ is the number of pinched geodesics disconnecting the surface and $\ell_{1}(g) \leqslant \ell_{2}(g) \leqslant \ldots \leqslant \ell_{n_{s}}(g)$ are the lengths of these separating geodesics ranked by increasing order.

Below, we take the convention that we repeat each eigenvalue according to its multiplicity, thus $\lambda_{j}(g)$ can be equal to $\lambda_{j+1}(g)$, and similarly for the $\nu_{j}(g)$.

Lemma 2.5. For each $g \in U_{g_{0}}$, let $v_{0}=(4 \pi(\mathbf{g}-1))^{-1}$ and $v_{1}, \ldots, v_{m} \in \mathbb{R}^{m+1}$ so that $\left(v_{i}\right)_{i=0, \ldots, m}$ is an orthonormal basis of eigenvectors for $Q_{g}$ with $v_{i}$ associated to $\nu_{i}(g)$. There is $C>0$ and $L \in \mathbb{N}$ such that for all $g \in U_{g_{0}}$, there exists an orthonormal basis $\varphi_{1}, \ldots, \varphi_{m}$ of $\oplus_{j=1}^{m} \operatorname{ker}\left(\Delta_{g}-\lambda_{j}(g)\right)$ satisfying

$$
\left\|f_{v_{j}}-\varphi_{j}\right\|_{L^{2}(M, g)} \leqslant C \delta^{\frac{1}{L}}, \text { and } \sum_{\lambda_{j}(g) \leqslant \delta} \frac{\Pi_{\lambda_{j}(g)}\left(x, x^{\prime}\right)}{\lambda_{j}(g)}=\sum_{j=1}^{m} \frac{f_{v_{j}}(x) f_{v_{j}}\left(x^{\prime}\right)}{\nu_{j}(g)}+\mathcal{O}\left(\frac{\delta^{\frac{1}{L}}}{\nu_{1}(g)}\right)
$$

where the error term is in $L^{\infty}$ norm on compact sets disjoints from $\cup_{j} \mathcal{C}_{j}(g)$.
Proof. To simplify notation we denote by $f_{j}$ the function $f_{v_{j}}$. We construct the basis $\varphi_{j}$ by an inductive process. Let $\left(\phi_{j}\right)_{j \in \mathbb{N}_{0}}$ be an orthonormal basis of $L^{2}(M, g)$ of eigenvectors for $\Delta_{g}$, ie.
$\Delta_{g} \phi_{j}=\lambda_{j}(g) \phi_{j}$. By Lemma 2.4, we have for $k \leqslant m$

$$
\left\|d f_{k}\right\|_{L^{2}}^{2}=\sum_{j=1}^{\infty} \lambda_{j}(g)\left\langle f_{k}, \phi_{j}\right\rangle^{2}=Q_{g}\left(v_{k}\right)(1+\mathcal{O}(\delta|\log \delta|))=\nu_{k}(g)\left\|f_{k}\right\|_{L^{2}}^{2}(1+\mathcal{O}(\delta|\log \delta|))
$$

and by Theorem 2.3, this gives for each $k=1, \ldots, m$

$$
\begin{equation*}
\sum_{j=1}^{\infty}\left(\frac{\lambda_{j}(g)}{\lambda_{k}(g)}-1\right)\left\langle f_{k}, \phi_{j}\right\rangle^{2}=\mathcal{O}\left(\delta^{\frac{1}{2}}\left\|f_{k}\right\|_{L^{2}}^{2}\right) \tag{2.32}
\end{equation*}
$$

If $\left|\frac{\lambda_{2}}{\lambda_{1}}-1\right|>\delta^{\frac{1}{4}}$, we set $\varphi_{1}:=\phi_{1}$ and $i_{1}=1$. By (2.32) with $k=1$, we get $\sum_{j=2}^{\infty}\left\langle f_{k}, \phi_{j}\right\rangle^{2}=\mathcal{O}\left(\delta^{\frac{1}{4}}\right)$ and thus $f_{1}= \pm \phi_{1}+\mathcal{O}_{L^{2}}\left(\delta^{\frac{1}{8}}\right)$. Since $\left\langle f_{i}, f_{j}\right\rangle_{L^{2}}=\mathcal{O}(\delta|\log \delta|)$ for $i \neq j$ by Lemma 2.4, we get $\left\langle f_{i}, \varphi_{1}\right\rangle=\mathcal{O}\left(\delta^{\frac{1}{8}}\right)$ for all $i>1$. If $\left|\frac{\lambda_{2}}{\lambda_{1}}-1\right| \leqslant \delta^{\frac{1}{4}}$, we let $i_{1} \geqslant 2$ be the smallest integer such that for each $i \leqslant i_{1},\left|\frac{\lambda_{i}}{\lambda_{i-1}}-1\right| \leqslant \delta^{\frac{1}{2^{i}}}$ and $\left|\frac{\lambda_{i_{1}+1}}{\lambda_{i_{1}}}-1\right|>\delta^{\frac{1}{2^{i_{1}+1}}}$, clearly $i_{1} \leqslant m$ since there are $m$ small eigenvalues. We define

$$
\varphi_{1}=\frac{\sum_{j=1}^{i_{1}}\left\langle f_{1}, \phi_{j}\right\rangle \phi_{j}}{\left\|\sum_{j=1}^{i_{1}}\left\langle f_{1}, \phi_{j}\right\rangle \phi_{j}\right\|}, \text { and } \widetilde{\varphi}_{i}=\frac{\sum_{j=1}^{i_{1}}\left\langle f_{i}, \phi_{j}\right\rangle \phi_{j}}{\left\|\sum_{j=1}^{i_{1}}\left\langle f_{i}, \phi_{j}\right\rangle \phi_{j}\right\|} \text { for } 1 \leqslant i \leqslant i_{1}
$$

Then we construct $\varphi_{2}, \ldots, \varphi_{i_{1}}$ by the Gram-Schmidt orthonormalization process from $\widetilde{\varphi}_{2}, \ldots, \widetilde{\varphi}_{i_{1}}$. Since $\left|\frac{\lambda_{i_{1}+1}}{\lambda_{i_{1}}}-1\right|>\delta^{\frac{1}{2_{1}+1}}$ and $\left|\frac{\lambda_{i_{1}}}{\lambda_{1}}-1\right|=\mathcal{O}\left(\delta{ }^{\frac{1}{2^{i_{1}}}}\right),(2.32)$ tells us that for each $i \leqslant i_{1}$,

$$
f_{i}=\sum_{j=1}^{i_{1}}\left\langle f_{i}, \phi_{j}\right\rangle \phi_{j}+\mathcal{O}_{L^{2}}\left(\delta^{\frac{1}{2_{1}+2}}\right)
$$

and thus $\widetilde{\varphi}_{i}=f_{i}+\mathcal{O}_{L^{2}}\left(\delta_{2^{\frac{1}{i_{1}+2}}}\right)$ for $i=1, \ldots, i_{1}$. Since $\left\langle f_{i}, f_{j}\right\rangle_{L^{2}}=\delta_{i j}+\mathcal{O}(\delta|\log \delta|)$ by Lemma 2.4, we deduce that for $i=1, \ldots, i_{1}$

$$
\varphi_{i}=f_{i}+\mathcal{O}_{L^{2}}\left(\delta_{2^{\frac{1}{i_{1}+2}}}\right)
$$

Now we prove the induction process in a way similar to the first step. Suppose we have constructed an orthornormal basis $\varphi_{1}, \ldots, \varphi_{\ell}$ of $\oplus_{j=1}^{\ell} \mathbb{R} \phi_{j}$ so that $\varphi_{j}=f_{j}+\mathcal{O}_{L^{2}}\left(\delta^{\frac{1}{L}}\right)$ for some $L \in \mathbb{N}$ and $\ell<m$. Notice that $\left\langle f_{k}, \phi_{j}\right\rangle=\mathcal{O}\left(\delta \frac{1}{L}\right)$ for all $k \geqslant \ell+1$ and $j \leqslant \ell$ by the induction assumption. Then if $\left|\frac{\lambda_{\ell+1}}{\lambda_{\ell}}-1\right|>\delta^{\frac{1}{L}}$, (2.32) with $k=\ell+1$ gives $\sum_{j=\ell+2}^{\infty}\left\langle f_{\ell+1}, \phi_{j}\right\rangle^{2}=\mathcal{O}\left(\delta^{\frac{1}{L}}\right)$, thus if we set $i_{\ell+1}=\ell+1$ and

$$
\varphi_{\ell+1}=\frac{\phi_{\ell+1}-\sum_{j=1}^{\ell}\left\langle\phi_{\ell+1}, \varphi_{j}\right\rangle \varphi_{j}}{\left\|\phi_{\ell+1}-\sum_{j=1}^{\ell}\left\langle\phi_{\ell+1}, \varphi_{j}\right\rangle \varphi_{j}\right\|}
$$

we get $\varphi_{\ell+1}=f_{\ell+1}+\mathcal{O}_{L^{2}}\left(\delta^{\frac{1}{2 L}}\right)$ and we have increased the induction step by 1. If $\left|\frac{\lambda_{\ell+1}}{\lambda_{\ell}}-1\right| \leqslant \delta^{\frac{1}{L}}$, we let $i_{\ell+1} \leqslant m$ be the smallest integer such that for all $i=\ell+1, \ldots, i_{\ell+1},\left|\frac{\lambda_{i}}{\lambda_{i-1}}-1\right| \leqslant \delta^{\frac{1}{L 2^{i-\ell-1}}}$ and $\left|\frac{\lambda_{i_{\ell+1}+1}}{\lambda_{i_{\ell+1}}}-1\right|>\delta \frac{1}{L_{2^{2} \ell+1^{-}-\ell}}$, and we will construct $\varphi_{\ell+1}, \ldots, \varphi_{i_{\ell+1}}$. Let $L^{\prime}=L 2^{i_{\ell+1}-\ell}$ and define

$$
\varphi_{\ell+1}=\frac{\sum_{j=\ell+1}^{i_{\ell+1}}\left\langle f_{\ell+1}, \phi_{j}\right\rangle \phi_{j}}{\left\|\sum_{j=\ell+1}^{i_{\ell+1}}\left\langle f_{\ell+1}, \phi_{j}\right\rangle \phi_{j}\right\|}, \text { and } \widetilde{\varphi}_{i}=\frac{\sum_{j=\ell+1}^{i_{\ell+1}}\left\langle f_{i}, \phi_{j}\right\rangle \phi_{j}}{\left\|\sum_{j=\ell+1}^{i_{\ell+1}}\left\langle f_{i}, \phi_{j}\right\rangle \phi_{j}\right\|} \text { for } \ell+1 \leqslant i \leqslant i_{\ell+1} .
$$

Then we construct $\varphi_{\ell+2}, \ldots, \varphi_{i_{\ell+1}}$ by the Gram-Schmidt orthonormalization process from $\widetilde{\varphi}_{\ell+2}, \ldots, \widetilde{\varphi}_{i_{\ell+1}}$. By induction assumption and $\left|\frac{\lambda_{i_{\ell+1}+1}}{\lambda_{i_{\ell+1}}}-1\right|>\delta^{\frac{1}{L^{\prime}}},(2.32)$ tells us that for each $i=\ell+1, \ldots, i_{\ell+1}$,

$$
f_{i}=\sum_{j=\ell+1}^{i_{\ell+1}}\left\langle f_{i}, \phi_{j}\right\rangle \phi_{j}+\mathcal{O}_{L^{2}}\left(\delta^{\frac{1}{2 L^{\prime}}}\right)
$$

and thus $\widetilde{\varphi}_{i}=f_{i}+\mathcal{O}_{L^{2}}\left(\delta^{\frac{1}{2 L^{\prime}}}\right)$ for $i=\ell+1, \ldots, i_{\ell+1}$. Since $\left\langle f_{i}, f_{j}\right\rangle_{L^{2}}=\delta_{i j}+\mathcal{O}(\delta|\log \delta|)$ by Lemma 2.4, we deduce that for $i=\ell+1, \ldots, i_{\ell+1}$

$$
\varphi_{i}=f_{i}+\mathcal{O}_{L^{2}}\left(\delta^{\frac{1}{2 L^{\prime}}}\right)
$$

and we have increased the induction step by $i_{\ell+1}-(\ell+1) \geqslant 1$. This inductive construction produces a sequence of integers $j_{0}=1, j_{1}=i_{1}, j_{2}=i_{i_{1}+1}, \ldots, j_{N}=m$ and $N$ associated blocks $E_{1}, \ldots E_{N}$, with $E_{k}=\left\{\varphi_{j_{k}}, \ldots, \varphi_{j_{k+1}}\right\}$ where the span of elements in $E_{k}$ is the span of $\left\{\phi_{j_{k}}, \ldots, \phi_{j_{k+1}}\right\}$. By construction we have

$$
\begin{aligned}
\sum_{\lambda_{j}(g) \leqslant \delta} \frac{\Pi_{\lambda_{j}}\left(x, x^{\prime}\right)}{\lambda_{j}(g)}= & \sum_{k=0}^{N} \sum_{j=j_{k}}^{j_{k+1}} \frac{\varphi_{j}(x) \varphi_{j}\left(x^{\prime}\right)}{\lambda_{j}(g)}+\mathcal{O}\left(\delta^{1 / L}\right) \\
& =\sum_{k=0}^{N} \sum_{j=j_{k}}^{j_{k+1}} \frac{\varphi_{j}(x) \varphi_{j}\left(x^{\prime}\right)}{\nu_{j}(g)}+\mathcal{O}\left(\frac{\delta^{\frac{1}{L}}}{\nu_{1}(g)}\right) \\
& =\sum_{j=1}^{m} \frac{f_{j}(x) f_{j}\left(x^{\prime}\right)}{\nu_{j}(g)}+\mathcal{O}\left(\frac{\delta^{\frac{1}{L}}}{\nu_{1}(g)}\right)
\end{aligned}
$$

for some $L \in \mathbb{N}$ large. Here we notice that the $\mathcal{O}\left(\frac{\delta^{\frac{1}{L}}}{\nu_{1}(g)}\right)$ can be taken in $L^{\infty}$ norm since $L^{2}$ norms on eigenfunctions give directly uniform $L^{\infty}$ norms on compact sets outside the collars $\mathcal{C}_{j}(g)$.
2.6. Example: the case of genus 2. For pedagogical purposes only, let us discuss more particularly the case of genus $\mathbf{g}=2$. In this case there can be only 3 simple curves in a partition and the maximal number of connected components separated by these curves is 2 : either 1 curve separates $M$ into two surfaces of genus 1 with 1 boundary component (Case 1) or two hyperbolic pants with 3 boundary components (Case 2). Consequently, the number of eigenvalues approaching 0 when we approach $\partial \overline{\mathcal{M}}_{2}$ is $m \in\{1,2\}$ (including the eigenvalue $\lambda=0$ ), we call them $\lambda_{0}=0$ and $\lambda_{1}(g)>0$ when $m=2$.

In Case 1 , take any partition $\mathrm{SP}_{1}$ by $\gamma_{1}, \gamma_{2}, \gamma_{3}$ with $\gamma_{1}$ being the only separating curve, and denote by $\ell_{j}(g)$ the length of the geodesic for $g$ freely homotopic to $\gamma_{j}$. We have

$$
\begin{equation*}
\lambda_{1}(g) \sim c_{1} \ell_{1}(g), \text { as } \ell_{1}(g) \rightarrow 0 \text { with } g \in U\left(\mathrm{SP}_{1}\right) \tag{2.33}
\end{equation*}
$$

where $c_{1}>0$ depending only on $\mathbf{g}=2$.
In Case 2, take $\mathrm{SP}_{2}$ any partition where $\gamma_{1}, \gamma_{2}, \gamma_{3}$ are separating simple curves and, if $\ell_{j}(g)$ is the length of the geodesic for $g$ homotopic to $\gamma_{j}$, then by [Bu1],

$$
\begin{equation*}
\lambda_{1}(g) \sim c_{2}\left(\ell_{1}(g)+\ell_{2}(g)+\ell_{3}(g)\right), \text { as }\left(\ell_{1}(g), \ell_{2}(g), \ell_{3}(g)\right) \rightarrow 0 \text { with } g \in U\left(\mathrm{SP}_{2}\right) \tag{2.34}
\end{equation*}
$$



Figure 1. On the left: Case 1. On the right: Case 2
for some $c_{2}>0$ depending only on $\mathbf{g}=2$.
In both cases, if $g \rightarrow g_{0}$ with $\left(M_{0}, g_{0}\right)$ a surface with nodes, $\left(M_{0}, g_{0}\right)$ decomposes into two finite volume hyperbolic surfaces $S_{1}$ and $S_{2}$, with volume $2 \pi$ (by Gauss-Bonnet), and by Proposition 2.2

$$
\begin{equation*}
\Pi_{\lambda_{1}(g)}\left(x, x^{\prime}\right) \rightarrow \frac{1}{4 \pi}\left(\mathbb{1}_{S_{1}}(x)-\mathbb{1}_{S_{2}}(x)\right)\left(\mathbb{1}_{S_{1}}\left(x^{\prime}\right)-\mathbb{1}_{S_{2}}\left(x^{\prime}\right)\right) \text { as } g \rightarrow g_{0} \tag{2.35}
\end{equation*}
$$

uniformly in ( $x, x^{\prime}$ ) on compact sets of $M_{0} \times M_{0}$.
2.7. Green's function near pinched geodesics. For later purpose, we will need a more detailed description of the Green's function $G_{g}$ than Proposition 2.2, in particular we shall need to know the behaviour of $G_{g}$ near the pinched geodesics. A recent work of Albin-RochonSher [ARS] gives a parametrix for $\Delta_{g}-\lambda$ when $\lambda$ is outside the spectrum and when there is one pinched geodesic (or the geodesics are pinched at the same speed). Here we need to know the behaviour in all possible directions of approach of the boundary of $\mathcal{M}_{\mathbf{g}}$. Our estimates are done in a way to be applied later for the study of the Gaussian multiplicative chaos measure in the cusp. The work of Melrose-Zhu [MeZh] gives a parametrix for some different Green function.

Let $\left(M_{0}, g_{0}\right)$ be a surface with nodes viewed as a surface with pairs of hyperbolic cusps, and let $\bar{U}_{g_{0}}$ be a local neighborhood of $g_{0}$ in $\overline{\mathcal{M}_{\mathrm{g}}}$ made of hyperbolic metrics $g_{s, t}$ as explained in Section 2.2, and denote $U_{g_{0}}=\mathcal{M}_{\mathbf{g}} \cap \bar{U}_{g_{0}}$. For convenience, we remove the parameters $(s, t)$ and just write $g$ for $g_{s, t}$. As in Proposition 2.2, we set

$$
\begin{equation*}
A_{g}\left(x, x^{\prime}\right):=G_{g}\left(x, x^{\prime}\right)-\sum_{j=1}^{m} \frac{\Pi_{\lambda_{j}}\left(x, x^{\prime}\right)}{\lambda_{j}}=\frac{1}{2 \pi i} \int_{\partial D(0, \varepsilon)} \frac{R_{g}\left(\lambda ; x, x^{\prime}\right)}{\lambda} d \lambda \tag{2.36}
\end{equation*}
$$

where $\lambda_{j}=\lambda_{j}(g)$ are the small eigenvalues tending to 0 as $g$ approaches the boundary of moduli space, $D(0, \varepsilon)$ is a small disk containing these eigenvalues (and only these ones) and $R_{g}\left(\lambda ; x, x^{\prime}\right)$ is the integral kernel of the resolvent of $\Delta_{g}$.

The hyperbolic surface $(M, g)$ decomposes into $M=S(g) \cup_{i=1}^{n_{p}} \mathcal{C}_{j}(g)$ where $\mathcal{C}_{j}(g)$ are the collars isometric to $[-1,1]_{\rho} \times(\mathbb{R} / \mathbb{Z})_{\theta}$ close to a given curve $\gamma_{j}$, where the metric $g$ is given (in geodesic normal coordinates to the geodesic $\gamma_{j}(g)$ homotopic to $\gamma_{j}$ ) by

$$
\begin{equation*}
g_{j}=\frac{d \rho^{2}}{\rho^{2}+\ell_{j}^{2}}+\left(\rho^{2}+\ell_{j}^{2}\right) d \theta^{2}, \tag{2.37}
\end{equation*}
$$

where $\ell_{j}=\ell_{j}(g)$ is the length of the geodesic $\gamma_{j}(g)$ and $S(g)$ is a compact manifold with boundary contained in a fixed (independent of $g$ ) compact set of $M_{0}$. This is also isometric to (by the coordinates change $\rho=\ell_{j} \sinh (t)$ )

$$
\left[-d_{j}, d_{j}\right]_{t} \times(\mathbb{R} / \mathbb{Z})_{\theta}, \quad g_{j}=d t^{2}+\ell_{j}^{2} \cosh ^{2}(t) d \theta^{2}, \quad \sinh \left(d_{j}\right)=1 / \ell_{j}
$$

The complete hyperbolic cylinder $\left\langle z \mapsto e^{\ell_{j}} z\right\rangle \backslash \mathbb{H}^{2}$ is isometric to

$$
\mathcal{F}_{j}:=\left(\mathbb{R}_{\rho} \times(\mathbb{R} / \mathbb{Z})_{\theta}, \quad g_{j}=\frac{d \rho^{2}}{\rho^{2}+\ell_{j}^{2}}+\left(\rho^{2}+\ell_{j}^{2}\right) d \theta^{2}\right)
$$

Notice that as $\ell_{j} \rightarrow 0$, the Riemannian manifold $\mathcal{F}_{j} \backslash\{\rho=0\}$ converges smoothly to two disconnected surfaces $(0, \infty)_{\rho} \times(\mathbb{R} / \mathbb{Z})_{\theta}$ with metric $d \rho^{2} / \rho^{2}+\rho^{2} d \theta^{2}$, which are isometric to two disjoint elementary quotients $\mathbb{H}^{2} /\langle z \mapsto z+1\rangle$. The Laplacian $\Delta_{g_{j}}$ is self-adjoint with spectrum $[1 / 4, \infty)$ (its self-adjoint realisation is through Friedrichs extension on $C_{c}^{\infty}\left(\mathcal{F}_{j}\right)$ ). It is invertible on $L^{2}\left(\mathcal{F}_{j}\right)$ and we can consider its resolvent $R_{g_{j}}(\lambda): L^{2}\left(\mathcal{F}_{j}\right) \rightarrow L^{2}\left(\mathcal{F}_{j}\right)$ which is holomorphic for $\lambda \notin[1 / 4, \infty)$. This is studied in details for example in [Bo, Prop. 5.2] or [GuZw, Appendix]: writing $\lambda=s(1-s)$ for $s$ close to 1 , we have

$$
R_{g_{j}}\left(\lambda ; \rho, \theta, \rho^{\prime}, \theta^{\prime}\right)=\sum_{k \in \mathbb{Z}} u_{k}\left(s ; \rho, \rho^{\prime}\right) e^{2 \pi i k\left(\theta-\theta^{\prime}\right)}
$$

for some explicit functions $u_{k}$ analytic in $s$ for $s$ close to 1 . We denote by $G_{g_{j}}$ the Green function corresponding to $\lambda=0$ (i.e. $s=1$ ). We give a more explicit bound at $\lambda=0$ in the following

Lemma 2.6. For $\ell_{j} \leqslant \min \left(|\rho|,\left|\rho^{\prime}\right|\right) \leqslant \max \left(|\rho|,\left|\rho^{\prime}\right|\right) \leqslant 1$ with $\rho \rho^{\prime}>0$, the Green function $G_{g_{j}}$ for the cylinder satisfies

$$
\begin{align*}
G_{g_{j}}\left(\rho, \theta, \rho^{\prime}, \theta^{\prime}\right)= & -\frac{1}{2 \pi} \log \left|1-e^{-\frac{2 \pi}{\ell_{j}}\left|\arctan \left(\frac{\rho}{\ell_{j}}\right)-\arctan \left(\frac{\rho^{\prime}}{\ell_{j}}\right)\right|+2 \pi i\left(\theta-\theta^{\prime}\right)}\right|  \tag{2.38}\\
& +\frac{1}{\ell_{j}} \min \left(F\left(\frac{|\rho|}{\ell_{j}}\right), F\left(\frac{\left|\rho^{\prime}\right|}{\ell_{j}}\right)\right)-\frac{1}{\pi \ell_{j}} F\left(\frac{|\rho|}{\ell_{j}}\right) F\left(\frac{\left|\rho^{\prime}\right|}{\ell_{j}}\right)+\mathcal{O}(1)
\end{align*}
$$

where $F(x):=\int_{x}^{\infty} \frac{d u}{1+u^{2}}$ and the remainder is uniform. If $\ell_{j} \leqslant \min \left(|\rho|,\left|\rho^{\prime}\right|\right) \leqslant \max \left(|\rho|,\left|\rho^{\prime}\right|\right) \leqslant 1$ with $\rho^{\prime} \rho<0$, then

$$
\begin{equation*}
G_{g_{j}}\left(\rho, \theta, \rho^{\prime}, \theta^{\prime}\right)=\frac{1}{\pi \ell_{j}} F\left(\frac{|\rho|}{\ell_{j}}\right) F\left(\frac{\left|\rho^{\prime}\right|}{\ell_{j}}\right)+\mathcal{O}\left(e^{-\frac{\pi}{2 \ell_{j}}}\right) \tag{2.39}
\end{equation*}
$$

If $|\rho| \in[1 / 2,1],\left|\rho^{\prime}\right| \leqslant 1$ and $\left|\rho-\rho^{\prime}\right|>\delta$ for some $\delta>0$, then there is $C$ depending only on $\delta$ so that

$$
\begin{equation*}
\left|\partial_{\rho} G_{g_{j}}\left(\rho, \theta, \rho^{\prime}, \theta^{\prime}\right)\right| \leqslant C \tag{2.40}
\end{equation*}
$$

If $\chi \in C_{c}^{\infty}(-1,1)$ and $\lambda \in[0,1)$, then we have the pointwise estimate

$$
\begin{equation*}
\int_{\mathbb{R} / \mathbb{Z}} \int_{-1}^{1} G_{g_{j}}\left(\rho, \theta, \rho^{\prime}, \theta^{\prime}\right) \chi\left(\rho^{\prime}\right) \rho^{\prime-\lambda} d \rho^{\prime} d \theta^{\prime} \leqslant 2\|\chi\|_{L^{\infty}}\left(\frac{\rho^{-\lambda}}{1-\lambda}+\frac{\rho^{-\lambda}-1}{\lambda}\right) \tag{2.41}
\end{equation*}
$$

where by convention $\frac{\rho^{-\lambda}-1}{\lambda}=|\log | \rho| |$ when $\lambda=0$. Finally, the Robin mass of $G_{g_{j}}$ satisfies

$$
\begin{equation*}
\left|m_{g_{j}}(\rho, \theta)-\frac{1}{\pi \ell_{j}}\left(\frac{\pi^{2}}{4}-\arctan \left(\frac{\rho}{\ell_{j}}\right)^{2}\right)-\frac{1}{2 \pi} \log \left(\sqrt{\rho^{2}+\ell_{j}^{2}}\right)\right| \leqslant C . \tag{2.42}
\end{equation*}
$$

Proof. Let $L$ be the operator $L=\left(\rho^{2}+\ell_{j}^{2}\right) \Delta_{g_{j}}$ acting on $T:=\mathbb{R}_{\rho} \times(\mathbb{R} / \mathbb{Z})_{\theta}$ with the measure $\left(\rho^{2}+\ell_{j}^{2}\right)^{-1} d \rho d \theta$, it is symmetric on $C_{c}^{\infty}(T)$. Changing coordinates to $t=\ell_{j}^{-1} \arctan \left(\rho / \ell_{j}\right), L$ becomes the operator

$$
L=-\partial_{t}^{2}-\partial_{\theta}^{2} \text { on }\left(-\frac{\pi}{2 \ell_{j}}, \frac{\pi}{2 \ell_{j}}\right)_{t} \times(\mathbb{R} / \mathbb{Z})_{\theta}
$$

with the measure $d t d \theta$. It is not self-adjoint but we can we consider the Friedrichs self-adjoint extension, which amounts to set Dirichlet conditions at $t= \pm \pi / 2 \ell_{j}$. It is clearly invertible for each $\ell_{j}>0$ and the inverse can be computed using Fourier decomposition in $\theta$. If $L^{-1}$ is written under the form

$$
\left(L^{-1} f\right)(t, \theta)=\int_{-\frac{\pi}{2 \ell_{j}}}^{\frac{\pi}{2 \ell_{j}}} \int_{\mathbb{R} / \mathbb{Z}} G_{L}\left(t, \theta, t^{\prime}, \theta^{\prime}\right) f\left(t^{\prime}, \theta^{\prime}\right) d \theta^{\prime} d t^{\prime}
$$

for some Green kernel $G_{L}$, then $G_{g_{j}}$ can be written as

$$
G_{g_{j}}\left(\rho, \theta, \rho^{\prime}, \theta^{\prime}\right)=G_{L}\left(\frac{\arctan \left(\frac{\rho}{\ell_{j}}\right)}{\ell_{j}}, \theta, \frac{\arctan \left(\frac{\rho^{\prime}}{\ell_{j}}\right)}{\ell_{j}}, \theta^{\prime}\right)
$$

This is clear since the left-hand side maps $C_{c}^{\infty}\left(\mathcal{F}_{j}\right)$ to $L^{2}\left(\mathcal{F}_{j}\right)$ and is a right inverse for $\Delta_{g_{j}}$ on $C_{c}^{\infty}\left(\mathcal{F}_{j}\right)$. Now, computing $G_{L}$ is quite simple: using the Fourier decomposition

$$
L^{-1} f(t, \theta)=\sum_{k \in \mathbb{Z}} e^{2 \pi i k \theta}\left(L_{k}^{-1} f_{k}\right)(t)
$$

where $f(t, \theta)=\sum_{k} e^{2 \pi i k \theta} f_{k}(t)$ and $L_{k}$ is the operator on $I_{j}:=\left(-\frac{\pi}{2 \ell_{j}}, \frac{\pi}{2 \ell_{j}}\right)$ given by $L_{k}=-\partial_{t}^{2}+$ $4 \pi^{2} k^{2}$ with Dirichlet condition at $\partial I_{j}$. For $k \neq 0$, straightforward Sturm-Liouville argument gives the expression of the Green function for $L_{k}$ : with $\bar{k}=2 \pi k$,

$$
\begin{aligned}
G_{L_{k}}\left(t, t^{\prime}\right)= & \frac{1}{2 \bar{k}\left(1-e^{-2 \bar{k} \pi / \ell_{j}}\right)}\left(\left(e^{-\bar{k} t}-e^{\bar{k}\left(t-\pi / \ell_{j}\right)}\right)\left(e^{\bar{k} t^{\prime}}-e^{-\bar{k}\left(t^{\prime}+\pi / \ell_{j}\right)}\right) \mathbb{1}_{t} \geqslant t^{\prime}\right. \\
& \left.+\left(e^{-\bar{k} t^{\prime}}-e^{\bar{k}\left(t^{\prime}-\pi / \ell_{j}\right)}\right)\left(e^{\bar{k} t}-e^{-\bar{k}\left(t+\pi / \ell_{j}\right)}\right) \mathbb{1}_{t^{\prime}} \geqslant t\right) \\
= & \frac{e^{-\bar{k}\left|t-t^{\prime}\right|}}{2 \bar{k}}+\frac{e^{-2 \pi \bar{k} / \ell_{j}} \cosh \left(\bar{k}\left(t-t^{\prime}\right)\right)-e^{-\bar{k} \pi / \ell_{j}} \cosh \left(\bar{k}\left(t+t^{\prime}\right)\right)}{\bar{k}\left(1-e^{-2 \pi \bar{k} / \ell_{j}}\right)} .
\end{aligned}
$$

If $\max \left(|\rho|,\left|\rho^{\prime}\right|\right) \leqslant 1$, we have $|t| \leqslant \frac{\pi}{2 \ell_{j}}-1+\mathcal{O}\left(\ell_{j}^{2}\right)$ and same for $t^{\prime}$ thus for $\ell_{j}$ small,

$$
\left|\frac{e^{-\bar{k} \pi / \ell_{j}} \cosh \left(\bar{k}\left(t+t^{\prime}\right)\right)}{\bar{k}\left(1-e^{-2 \pi \bar{k} / \ell_{j}}\right)}\right| \leqslant \frac{e^{-\bar{k} / 2}}{\bar{k}}, \quad\left|\frac{e^{-2 \pi \bar{k} / \ell_{j}} \cosh \left(\bar{k}\left(t-t^{\prime}\right)\right)}{\bar{k}\left(1-e^{-2 \pi \bar{k} / \ell_{j}}\right)}\right| \leqslant \frac{e^{-\pi \bar{k} / \ell_{j}}}{\bar{k}} .
$$

We then get

$$
\sum_{k \neq 0} G_{L_{k}}\left(t, t^{\prime}\right) e^{i \bar{k}\left(\theta-\theta^{\prime}\right)}=\sum_{k \neq 0} \frac{e^{-\bar{k}\left|t-t^{\prime}\right|+i \bar{k}\left(\theta-\theta^{\prime}\right)}}{2 \bar{k}}+\mathcal{O}(1)=-\frac{1}{2 \pi} \log \left|1-e^{-2 \pi\left(\left|t-t^{\prime}\right|-i\left(\theta-\theta^{\prime}\right)\right)}\right|+\mathcal{O}(1)
$$

when $\ell_{j}$ is small. Notice also that if $|\rho| \in\left[\ell_{j}, 1\right]$ and $\left|\rho^{\prime}\right| \in\left[\ell_{j}, 1\right]$ with $\rho$ and $\rho^{\prime}$ having different sign, then

$$
\begin{equation*}
\left|\sum_{k \neq 0} G_{L_{k}}\left(t(\rho), t^{\prime}\left(\rho^{\prime}\right)\right) e^{i \bar{k}\left(\theta-\theta^{\prime}\right)}\right| \leqslant C e^{-\frac{\pi}{2 \ell_{j}}} \tag{2.43}
\end{equation*}
$$

if $t(\rho)=\ell_{j}^{-1} \arctan \left(\rho / \ell_{j}\right)$ and similarly for $t^{\prime}\left(\rho^{\prime}\right)$. Next for $k=0$, the Green function is given by the expression

$$
G_{L_{0}}\left(t, t^{\prime}\right)=-\frac{1}{2}\left|t-t^{\prime}\right|-\frac{\ell_{j}}{\pi} t t^{\prime}+\frac{\pi}{4 \ell_{j}},
$$

thus we get

$$
G_{L_{0}}\left(t(\rho), t^{\prime}\left(\rho^{\prime}\right)\right)=-\frac{1}{2 \ell_{j}}\left|\arctan \left(\frac{\rho}{\ell_{j}}\right)-\arctan \left(\frac{\rho^{\prime}}{\ell_{j}}\right)\right|+\frac{1}{\pi \ell_{j}}\left(\frac{\pi^{2}}{4}-\arctan \left(\frac{\rho}{\ell_{j}}\right) \arctan \left(\frac{\rho^{\prime}}{\ell_{j}}\right)\right) .
$$

If $F(x):=\int_{x}^{\infty} \frac{d u}{1+u^{2}}$, this can be rewritten as

$$
\begin{align*}
G_{L_{0}}\left(\rho, \rho^{\prime}\right) & =-\frac{1}{2 \ell_{j}}\left|F\left(\frac{\rho}{\ell_{j}}\right)-F\left(\frac{\rho^{\prime}}{\ell_{j}}\right)\right|+\frac{1}{2 \ell_{j}}\left(F\left(\frac{\rho}{\ell_{j}}\right)+F\left(\frac{\rho^{\prime}}{\ell_{j}}\right)\right)-\frac{1}{\pi \ell_{j}} F\left(\frac{\rho}{\ell_{j}}\right) F\left(\frac{\rho^{\prime}}{\ell_{j}}\right) \\
& = \begin{cases}\frac{1}{\ell_{j}} \min \left(F\left(\frac{\rho}{\ell_{j}}\right), F\left(\frac{\rho^{\prime}}{\ell_{\ell^{\prime}}}\right)\right)-\frac{1}{\pi \ell_{j}} F\left(\frac{\rho}{\ell_{j}}\right) F\left(\frac{\rho^{\prime}}{\ell_{j}}\right), & \text { if } \rho>0, \rho^{\prime}>0 \\
\frac{1}{\pi \ell_{j}} F\left(\frac{\rho}{\ell_{j}}\right) F\left(-\frac{\rho^{\prime}}{\ell_{j}}\right), & \text { if } \rho>0, \rho^{\prime}<0\end{cases} \tag{2.44}
\end{align*}
$$

from which (2.38) and (2.39) follow.
Now, we also see from the expressions of $G_{L_{k}}$ above that if $|\rho| \in[1 / 2,1]$ and $\left|\rho-\rho^{\prime}\right|>\delta$ for some fixed constant $\delta>0$, then

$$
\left|\partial_{\rho} G_{g_{j}}\left(\rho, \theta, \rho^{\prime}, \theta^{\prime}\right)\right| \leqslant C
$$

for some constant $C$ depending only on $\delta$ but not on $\ell_{j}$.
Next we prove (2.41). Let $F(x)=\int_{x}^{\infty} \frac{d u}{1+u^{2}}$ and let $c:=\|\chi\|_{L^{\infty}}$. For $\rho>0$, we write

$$
\begin{aligned}
\int G_{g_{j}}\left(\rho, \theta, \rho^{\prime}, \theta^{\prime}\right) \frac{\chi\left(\rho^{\prime}\right)}{\rho^{\prime \lambda}} d x^{\prime} & =\frac{1}{\ell_{j} \pi} F\left(\frac{\rho}{\ell_{j}}\right) \int_{-1}^{0} F\left(-\frac{\rho^{\prime}}{\ell_{j}}\right) \frac{\chi\left(\rho^{\prime}\right)}{\rho^{\prime \lambda}} d \rho^{\prime}+\int_{0}^{\infty} G_{g_{j}}\left(\rho, \theta, \rho^{\prime}, \theta^{\prime}\right) \frac{\chi\left(\rho^{\prime}\right)}{\rho^{\prime \lambda}} d \rho^{\prime} \\
& \leqslant \frac{c}{\pi \rho} \int_{-1}^{0} F\left(-\frac{\rho^{\prime}}{\ell_{j}}\right) \frac{1}{\rho^{\prime \lambda}} d \rho^{\prime}+\frac{c}{\rho} \int_{0}^{\rho} \frac{d \rho^{\prime}}{\rho^{\prime \lambda}}+c \int_{\rho}^{1} \frac{d \rho^{\prime}}{\rho^{\prime 1+\lambda}} \\
& \leqslant 2 c \frac{\rho^{-\lambda}}{1-\lambda}+2 c \frac{\rho^{-\lambda}-1}{\lambda}
\end{aligned}
$$

where we used (2.44) in the second line and

$$
\int_{-1}^{-\rho} F\left(-\frac{\rho^{\prime}}{\ell_{j}}\right) \frac{d \rho^{\prime}}{\rho^{\prime \lambda}} \leqslant \ell_{j} \int_{\rho}^{1} \frac{d \rho^{\prime}}{\rho^{\prime 1+\lambda}} \leqslant \ell_{j} \frac{\rho^{-\lambda}-1}{\lambda}, \quad \int_{-\rho}^{0} F\left(-\frac{\rho^{\prime}}{\ell_{j}}\right) \frac{d \rho^{\prime}}{\rho^{\prime \lambda}} \leqslant \frac{\pi}{2(1-\lambda)} \rho^{1-\lambda}
$$

for the third line. The same estimate works in the case $\rho<0$.
Using the estimates and expressions above, we also get as $\left|\left(\rho^{\prime}, \theta^{\prime}\right)-(\rho, \theta)\right|$ is small

$$
G_{L}\left(t, \theta, t^{\prime}, \theta^{\prime}\right)=-\frac{1}{2 \pi} \log \left(\sqrt{\left(t-t^{\prime}\right)^{2}+\left(\theta-\theta^{\prime}\right)^{2}}\right)-\frac{\ell_{j}}{\pi} t^{2}+\frac{\pi}{4 \ell_{j}}+\mathcal{O}(1)+\mathcal{O}\left(\left|t-t^{\prime}\right|+\left|\theta-\theta^{\prime}\right|\right)
$$

where $\mathcal{O}(1)$ is independent of all variables. We also have

$$
\log \left(d_{g_{j}}\left(\rho, \theta, \rho^{\prime}, \theta^{\prime}\right)\right)=\log \left(\sqrt{\left(t-t^{\prime}\right)^{2}+\left(\theta-\theta^{\prime}\right)^{2}}\right)+\log \left(\sqrt{\rho^{2}+\ell_{j}^{2}}\right)+\mathcal{O}\left(\left|t-t^{\prime}\right|+\left|\theta-\theta^{\prime}\right|\right)
$$

thus the Robin mass of $G_{g_{j}}$ satisfies (2.42).
Next we express $\int_{\partial D(0, \varepsilon)} R_{g}(\lambda) d \lambda / \lambda$ in terms of $G_{g_{j}}$. Let $\chi_{j}, \chi_{j}^{\prime} \in C^{\infty}(M)$ which are supported in $\mathcal{C}_{j}(g)$, depending only on the variable $\rho$ associated to the metric $g$, are equal to 1 in $|\rho| \leqslant 1 / 4$ and such that $\chi_{j}^{\prime}=1$ on a neighborhood of $\operatorname{supp}\left(\chi_{j}\right)$. We will use the diffeomorphisms
between $\mathcal{C}_{j}(g)$ and the subset $|\rho|<1$ of $\mathcal{F}_{j}$ to identify these sets (for notational simplicity we won't input these diffeos below). Using that $\Delta_{g}=\Delta_{g_{j}}$ in $\mathcal{C}_{j}(g)$, we have, with $\chi:=\sum_{j} \chi_{j}$

$$
\left(\Delta_{g}-\lambda\right) \sum_{j} \chi_{j}^{\prime} R_{g_{j}}(\lambda) \chi_{j}=\chi+\sum_{j}\left[\Delta_{g}, \chi_{j}^{\prime}\right] R_{g_{j}}(\lambda) \chi_{j}
$$

and thus

$$
\sum_{j} \chi_{j}^{\prime} R_{g_{j}}(\lambda) \chi_{j}=R_{g}(\lambda) \chi+R_{g}(\lambda) \sum_{j}\left[\Delta_{g}, \chi_{j}^{\prime}\right] R_{g_{j}}(\lambda) \chi_{j} .
$$

Similarly, we also have

$$
\sum_{j} \chi_{j} R_{g_{j}}(\lambda) \chi_{j}^{\prime}=\chi R_{g}(\lambda)+\sum_{j} \chi_{j} R_{g_{j}}(\lambda)\left[\chi_{j}^{\prime}, \Delta_{g}\right] R_{g}(\lambda)
$$

and therefore (using also $\chi_{j}^{\prime} \chi_{k}^{\prime}=0$ if $i \neq k$ )

$$
\begin{aligned}
\chi R_{g}(\lambda) \chi= & \sum_{j} \chi_{j} R_{g_{j}}(\lambda) \chi_{j}-\sum_{j} \chi_{j} R_{g_{j}}(\lambda) \chi_{j}^{\prime}\left[\Delta_{g}, \chi_{j}^{\prime}\right] R_{g_{j}}(\lambda) \chi_{j} \\
& +\sum_{j} \chi_{j} R_{g_{j}}(\lambda)\left[\Delta_{g}, \chi_{j}^{\prime}\right] R_{g}(\lambda) \sum_{i}\left[\chi_{i}^{\prime}, \Delta_{g}\right] R_{g_{i}}(\lambda) \chi_{i} \\
= & \sum_{j} \chi_{j} R_{g_{j}}(\lambda) \chi_{j}-K_{1}(\lambda)+K_{2}(\lambda) R_{g}(\lambda) K_{3}(\lambda)
\end{aligned}
$$

where $K_{1}(\lambda), K_{2}(\lambda), K_{3}(\lambda)$ are defined by the equation. Remark that, in the right hand side, only the last term has poles (first order) in $D(0, \varepsilon)$. Therefore, using Cauchy formula

$$
\begin{aligned}
\int_{\partial D(0, \varepsilon)} \frac{\chi R_{g}(\lambda) \chi}{2 \pi i \lambda} d \lambda= & \sum_{j} \chi_{j} R_{g_{j}}(0) \chi_{j}-K_{1}(0)+K_{2}(0) A_{g} K_{3}(0) \\
& -K_{2}^{\prime}(0) \Pi_{0} K_{3}(0)-K_{2}(0) \Pi_{0} K_{3}^{\prime}(0) \\
& -\sum_{k=1}^{m} \frac{\left(K_{2}\left(\lambda_{k}\right)-K_{2}(0)\right)}{\lambda_{k}} \Pi_{\lambda_{k}} K_{3}\left(\lambda_{k}\right)-K_{2}(0) \Pi_{\lambda_{k}} \frac{\left(K_{3}\left(\lambda_{k}\right)-K_{3}(0)\right)}{\lambda_{k}}
\end{aligned}
$$

where $K_{i}^{\prime}(0):=\left.\partial_{\lambda} K_{i}(\lambda)\right|_{\lambda=0}$ and $A_{g}:=\int_{\partial D(0, \varepsilon)} \frac{R_{g}(\lambda)}{2 \pi i \lambda} d \lambda$. Let us analyse those terms more carefully.

Proposition 2.7. Let $g_{0} \in \partial \mathcal{M}_{\mathrm{g}}$ be a hyperbolic surface with cusps in the boundary of moduli space. Then there is a neighborhood $U_{g_{0}}$ of $g_{0}$ in $\mathcal{M}_{\mathbf{g}}$ and $C>0$ such that for all $g \in U_{g_{0}}$ and $x, x^{\prime} \in \cup_{j} \mathcal{C}_{j}(g)$

$$
\begin{aligned}
\chi(x) \chi\left(x^{\prime}\right) G_{g}\left(x, x^{\prime}\right)= & \sum_{j} \chi_{j}(x) G_{g_{j}}\left(x, x^{\prime}\right) \chi_{j}\left(x^{\prime}\right)+Q_{g}\left(x, x^{\prime}\right)-H_{0}(x) J_{0}\left(x^{\prime}\right)-J_{0}(x) H_{0}\left(x^{\prime}\right) \\
& +\sum_{k=1}^{m}\left(\frac{H_{k}(x)-\lambda_{k} J_{k}(x)}{\sqrt{\lambda_{k}}}\right)\left(\frac{H_{k}\left(x^{\prime}\right)-\lambda_{k} J_{k}\left(x^{\prime}\right)}{\sqrt{\lambda_{k}}}\right)
\end{aligned}
$$

with $Q_{g} \in C^{\infty}(M \times M), H_{k}=\chi \varphi_{\lambda_{k}}, J_{k} \in C^{\infty}(M)$ satisfying

$$
\begin{gathered}
\left|Q_{g}\right|_{L^{\infty}} \leqslant C, \quad\left|H_{0}\right|_{L^{\infty}} \leqslant C, \quad\left|H_{k}(\rho, \theta)\right| \leqslant C|\rho|^{-\left(1-s_{k}\right)}, \\
\left|J_{0}(\rho, \theta)\right| \leqslant C|\log | \rho \|, \quad\left|J_{k}(\rho, \theta)\right| \leqslant C\left(|\rho|^{-\left(1-s_{k}\right)}+\frac{|\rho|^{-\left(1-s_{k}\right)}-1}{\left(1-s_{k}\right)}\right)
\end{gathered}
$$

with $s_{k}\left(1-s_{k}\right)=\lambda_{k}(g)>0$ the small eigenvalues of $\Delta_{g}$ converging to 0 as $g$ approaches $\partial \mathcal{M}_{\mathbf{g}}$ and $s_{k}=1-\lambda_{k}(g)+\mathcal{O}\left(\lambda_{k}(g)^{2}\right)$. Here $\rho=\ell_{j} \sinh (t)$ in $\mathcal{C}_{j}(g)$ with $t$ the signed distance to the geodesic $\gamma_{j}(g)$.

Proof. The integral kernel of $\sum_{j} \chi_{j} R_{g_{j}}(0) \chi_{j}$ is $\chi_{j}(x) \chi_{j}\left(x^{\prime}\right) G_{g_{j}}\left(x, x^{\prime}\right)$ in $\mathcal{C}_{j}(g)$ with respect to the measure $\mathrm{dv}_{g_{j}}$. This term has an explicit bound by using Lemma 2.6. The function $Q_{g}\left(x, x^{\prime}\right)$ will be chosen to be the integral kernel of $-K_{1}(0)+K_{2}(0) A_{g} K_{3}(0)$, let us show this is smooth and uniformly bounded. The integral kernel of $K_{1}(0)$ is of the form

$$
K_{1}\left(0 ; x, x^{\prime}\right)=\sum_{j} \chi_{j}(x) \int_{\mathcal{C}_{j}^{\prime}} G_{g_{j}}(x, y) \chi_{j}^{\prime}(y) P_{j}(y) G_{g_{j}}\left(y, x^{\prime}\right) \operatorname{dv}_{g_{j}}(y) \chi_{j}\left(x^{\prime}\right)
$$

in $\mathcal{C}_{j}(g)$ with respect to the measure $\mathrm{dv}_{g_{j}}$, where $P_{j}$ is a smooth differential operator of order 1 supported in $\operatorname{supp}\left(\nabla \chi_{j}\right)$ (thus far from the $\gamma_{j}(g)=\{\rho=0\}$ curve). Thus by (2.40) and (2.39), we directly get that for all $x, x^{\prime} \in \mathcal{C}_{j}(g)$ with $|\rho(x)|>\ell_{j}$ and $\left|\rho\left(x^{\prime}\right)\right|>\ell_{j}$

$$
\left|K_{1}\left(0 ; x, x^{\prime}\right)\right| \leqslant C, \quad \text { and if } \rho(x) \rho\left(x^{\prime}\right)<0,\left|K_{1}\left(0 ; x, x^{\prime}\right)\right| \leqslant C \ell_{j} .
$$

By Proposition 2.2, the norm $\left\|A_{g}\right\|_{L^{2}(W) \rightarrow L^{2}(W)}$ is uniformly bounded if $W:=\operatorname{supp}\left(\nabla \chi_{j}^{\prime}\right)$, and by the same argument as for $K_{1}(0), K_{2}(0)$ and $K_{3}(0)$ have smooth integral kernels that are uniformly bounded with respect to $\ell_{j}$, thus there is $C>0$ uniform so that for all $x, x^{\prime} \in M$

$$
\left|\left(K_{2}(0) A_{g} K_{3}(0)\right)\left(x, x^{\prime}\right)\right| \leqslant C .
$$

We can rewrite $K_{2}^{\prime}(0) \Pi_{0} K_{3}(0)$ by using that $\Pi_{0}=c$ with $c=1 / \operatorname{Vol}_{g}(M)$ :

$$
\begin{aligned}
\left(K_{2}^{\prime}(0) \Pi_{0} K_{3}(0)\right)\left(x, x^{\prime}\right) & =c \sum_{j} \chi_{j}(x)\left(\partial_{\lambda} R_{g_{j}}(0) \Delta_{g} \chi_{j}^{\prime}\right)(x) \chi\left(x^{\prime}\right) \\
& =c \sum_{j} \chi_{j}(x)\left(R_{g_{j}}(0) \chi_{j}^{\prime}\right)(x) \chi\left(x^{\prime}\right)=H_{0}\left(x^{\prime}\right) J_{0}(x)
\end{aligned}
$$

where we used $\partial_{\lambda} R_{g_{j}}(0) \Delta_{g_{j}}=R_{g_{j}}(0), J_{0}(x):=\sqrt{c} \sum_{j} \chi_{j}(x) \int_{\mathcal{C}_{j}} G_{g_{j}}\left(x, x^{\prime}\right) \chi_{j}\left(x^{\prime}\right) \mathrm{dv}_{g_{j}}\left(x^{\prime}\right)$ and $H_{0}:=\sqrt{c} \chi$ are smooth functions on $M$. Similarly, we get

$$
\left(K_{2}(0) \Pi_{0} K_{3}^{\prime}(0)\right)\left(x, x^{\prime}\right)=H_{0}(x) J_{0}\left(x^{\prime}\right) .
$$

Now the bound (2.41) gives that $\left|J_{0}(\rho, \theta)\right| \leqslant C|\log | \rho| |$ and $\left|H_{0}(x)\right| \leqslant C$ for some uniform $C>0$.

Using that $\left(\Delta_{g}-\lambda_{k}\right) \Pi_{\lambda_{k}}=\Pi_{\lambda_{k}}\left(\Delta_{g}-\lambda_{k}\right)=0$, we get

$$
\begin{aligned}
\frac{\left(K_{2}\left(\lambda_{k}\right)-K_{2}(0)\right)}{\lambda_{k}} \Pi_{\lambda_{k}} K_{3}\left(\lambda_{k}\right) & =\sum_{j} \chi_{j} R_{g_{j}}(0) \chi_{j}^{\prime} \Pi_{\lambda_{k}} \sum_{i}\left[\chi_{i}^{\prime}, \Delta_{g}\right] R_{g_{i}}\left(\lambda_{k}\right) \chi_{i} \\
& =\sum_{j} \chi_{j} R_{g_{j}}(0) \chi_{j}^{\prime} \Pi_{\lambda_{k}} \chi
\end{aligned}
$$

Similarly, we also get

$$
\begin{aligned}
K_{2}(0) \Pi_{\lambda_{k}} \frac{\left(K_{3}\left(\lambda_{k}\right)-K_{3}(0)\right)}{\lambda_{k}} & =\sum_{j} \chi_{j} R_{g_{j}}(0)\left[\Delta_{g}, \chi_{j}^{\prime}\right] \Pi_{\lambda_{k}} \sum_{i} \chi_{i}^{\prime} R_{g_{i}}(0) \chi_{i} \\
& =\chi \Pi_{\lambda_{1}} \sum_{i} \chi_{i}^{\prime} R_{g_{i}^{\prime}}(0) \chi_{i}-\lambda_{k} \sum_{j} \chi_{j} R_{g_{j}}(0) \chi_{j}^{\prime} \Pi_{\lambda_{k}} \sum_{i} \chi_{i}^{\prime} R_{g_{i}}(0) \chi_{i} .
\end{aligned}
$$

Write $\Pi_{\lambda_{k}}\left(x, x^{\prime}\right)=\varphi_{\lambda_{k}}(x) \varphi_{\lambda_{k}}\left(x^{\prime}\right)$ for some $\varphi_{\lambda_{k}} \in \operatorname{ker}\left(\Delta_{g}-\lambda_{k}\right)$ with $L^{2}$ norm 1, then we get

$$
\begin{gathered}
\frac{\left(K_{2}\left(\lambda_{k}\right)-K_{2}(0)\right)}{\lambda_{k}} \Pi_{\lambda_{k}} K_{3}\left(\lambda_{k}\right)+K_{2}(0) \Pi_{\lambda_{k}} \frac{\left(K_{3}\left(\lambda_{k}\right)-K_{3}(0)\right)}{\lambda_{k}}= \\
H_{k}(x) J_{k}\left(x^{\prime}\right)+J_{k}(x) H_{k}\left(x^{\prime}\right)-\lambda_{k} J_{k}(x) J_{k}\left(x^{\prime}\right)
\end{gathered}
$$

where $H_{k}=\chi \varphi_{\lambda_{k}}, J_{k}(x)=\sum_{j} \chi_{j}(x) \int_{\mathcal{C}_{j}(g)} G_{g_{j}}\left(x, x^{\prime}\right) \chi_{j}^{\prime}(x) \varphi_{\lambda_{k}}\left(x^{\prime}\right) \mathrm{dv}_{g_{j}}\left(x^{\prime}\right)$. To conclude, we need some estimates on the eigenfunction $\varphi_{\lambda_{k}}$ associated to the small positive eigenvalue $\lambda_{k}$ of $\Delta_{g}$ on $M$

Lemma 2.8. Let $g_{0} \in \partial \mathcal{M}_{g}$ be a surface with node with $m+1$ connected components and denote by $\lambda_{1}, \ldots \lambda_{m}$ the positive small eigenvalues. For each $\varepsilon>0$, there is a neighborhood $U_{g_{0}}$ of $g_{0}$ and $C>0$ such that for each $g \in U_{g_{0}}$, the following holds: if $\varphi_{\lambda_{i}}$ is an eigenfunction for $\lambda_{i}$ which satisfies $\left|\varphi_{\lambda_{i}}\right|_{\rho= \pm 1}-a_{i j}^{ \pm} \mid<\varepsilon$ for some constants $a_{i j}^{ \pm} \in \mathbb{R}$ in the collar $\mathcal{C}_{j}(g)$, then it satisfies

$$
\varphi_{\lambda_{i}}(\rho, \theta)=\left(a_{i j}^{+} \mathbb{1}_{\rho>0}+a_{i j}^{-} \mathbb{1}_{\rho<0}\right)|\rho|^{s_{i}-1}(1+\mathcal{O}(\varepsilon))+\mathcal{O}(\varepsilon)
$$

in the region $\left\{|\rho|>C \ell_{j}\right\}$ of the collar $\mathcal{C}_{j}(g)$, where $s_{i}\left(1-s_{i}\right)=\lambda_{i}$ and $s_{i}=1-\lambda_{i}+\mathcal{O}\left(\lambda_{i}^{2}\right)$ when $\lambda_{i}$ is small. In the region $|\rho|<C \ell_{j}$, there is $C^{\prime}>0$ such that

$$
\left|\varphi_{\lambda_{i}}(\rho, \theta)\right| \leqslant C^{\prime}|\rho|^{s_{i}-1} .
$$

Proof. To simplify notations, we remove the $i, j$ indices from $a_{i j}^{ \pm}, \lambda_{i}$ and $s_{i}$. We decompose $\varphi_{\lambda}$ in Fourier modes in $\theta$ : there are $b_{k} \in C^{\infty}([-1,1])$ so that

$$
\varphi_{\lambda}(\rho, \theta)=\sum_{k=-\infty}^{\infty} b_{k}(\rho) e^{2 \pi i k \theta}
$$

and the series converges uniformly. Since $\varphi_{\lambda}$ can be supposed real-valued, $b_{0}$ is real and $b_{-k}=$ $\overline{b_{k}}$. Moreover $a_{k}(t):=b_{k}\left(\ell_{j} \sinh (t)\right)$ satisfies the ODE

$$
\left(-\partial_{t}^{2}-\tanh (t) \partial_{t}+\frac{4 \pi^{2} k^{2}}{\ell_{j}^{2} \cosh (t)^{2}}-\lambda\right) a_{k}(t)=0, \quad a_{k}\left( \pm t_{j}\right)=a_{k}^{ \pm}
$$

if $t_{j}$ is defined by $\ell_{j} \sinh \left(t_{j}\right)=1$. We write $a_{k}^{ \pm}:=\left.a_{k}\right|_{t= \pm t_{j}}=\left.b_{k}\right|_{\rho= \pm 1}$. First we make the following observation for each $k \neq 0$ :

$$
\begin{equation*}
\left|a_{k}(t)\right| \leqslant \max \left(\left|a_{k}^{+}\right|,\left|a_{k}^{-}\right|\right)<\varepsilon . \tag{2.45}
\end{equation*}
$$

Indeed, assume that $a_{k}$ achieves its maximum at $T \in\left(-t_{j}, t_{j}\right)$ with $a_{k}(T)>\max \left(a_{k}^{+}, a_{k}^{-}\right)$, then if $a_{k}(T)>0$

$$
-a_{k}^{\prime \prime}(T)=\left(\lambda-\frac{4 \pi^{2} k^{2}}{\ell_{j}^{2} \cosh (T)^{2}}\right) a_{k}(T)<0
$$

when $\ell_{j}$ is smaller than a uniform constant. This contradicts that $a_{k}(T)$ is a local maximum, thus $a_{k}$ achieves its maximum at $\pm t_{j}$ or its maximum is non-positive, in which case $a_{k}^{+} \leqslant 0$ and $a_{k}^{-} \leqslant 0$. In both cases, $\left|a_{k}(T)\right| \leqslant \max \left(\left|a_{k}^{+}\right|,\left|a_{k}^{-}\right|\right)$. The same argument works with the minimum and this shows (2.45). Next we analyze $u_{0}(t)$, and it is convenient for that to use $s \in(0,1)$ so that $s(1-s)=\lambda\left(\right.$ then $\left.s=1-\lambda+\mathcal{O}(\lambda)^{2}\right)$. There are 2 independent solutions of
the ODE with $k=0$ (and no bounday condition), the first one $v_{0}$ is odd in $t$, the other one $u_{0}$ is even, they are given by [Bo, Chapter 5.1]

$$
\begin{aligned}
v_{0}(t)= & \operatorname{sign}(t) \frac{\Gamma\left(\frac{1}{2}-s\right) \Gamma\left(\frac{1+s}{2}\right)^{2}}{\Gamma\left(s-\frac{1}{2}\right) \Gamma\left(1-\frac{s}{2}\right)^{2}}|\sinh (t)|^{-s} F\left(\frac{1+s}{2}, \frac{s}{2}, \frac{1}{2}+s ; \frac{-1}{\sinh (t)^{2}}\right) \\
& +\operatorname{sign}(t)|\sinh (t)|^{s-1} F\left(\frac{2-s}{2}, \frac{1-s}{2}, \frac{3}{2}-s ; \frac{-1}{\sinh (t)^{2}}\right), \\
u_{0}(t)= & \frac{\Gamma\left(\frac{1}{2}-s\right) \Gamma\left(\frac{s}{2}\right)^{2}}{\Gamma\left(s-\frac{1}{2}\right) \Gamma\left(\frac{1-s}{2}\right)^{2}}|\sinh (t)|^{-s} F\left(\frac{s}{2}, \frac{s+1}{2}, \frac{1}{2}+s ; \frac{-1}{\sinh (t)^{2}}\right) \\
& +|\sinh (t)|^{s-1} F\left(\frac{1-s}{2}, 1-\frac{s}{2}, \frac{3}{2}-s ; \frac{-1}{\sinh (t)^{2}}\right)
\end{aligned}
$$

where $F(a, b, c ; z)$ is the hypergeometric function, holomorphic in the variables $a, b, c$ for $a, b, c \in$ $\mathbb{C}$ in the half-plane $\mathbb{C}_{+}:=\{c \in \mathbb{C}, \operatorname{Re}(c)>0\}$ if $z \in(-\infty, 0)$, it is smooth in $z$ and for $z<0$ small

$$
F(a, b, c ; z)=1+\mathcal{O}(|z|)
$$

where the remainder is uniform for $a, b, c$ in compact sets of $\mathbb{C}_{+}$. In particular, there is $C>0$ uniform in $g$ so that for $|t|>C$,

$$
v_{0}(t)=\operatorname{sign}(t)|\sinh (t)|^{s-1}+\mathcal{O}\left(|\sinh t|^{-s}\right), \quad u_{0}(t)=|\sinh (t)|^{s-1}+\mathcal{O}\left(|\sinh t|^{-s}\right)
$$

where the remainder is uniform with respect to $\lambda$ for $\lambda>0$ small. We obtain

$$
a_{0}(t)=\frac{\left(a_{0}^{+}+a_{0}^{-}\right)}{2} \frac{u_{0}(t)}{u_{0}\left(t_{j}\right)}+\frac{\left(a_{0}^{+}-a_{0}^{-}\right)}{2} \frac{v_{0}(t)}{v_{0}\left(t_{j}\right)}
$$

and we deduce that

$$
a_{0}(t)=\left(a_{0}^{+} \mathbb{1}_{t>0}+a_{0}^{-} \mathbb{1}_{t<0}\right)\left|\ell_{j} \sinh (t)\right|^{s-1}+\mathcal{O}\left(\ell_{j}^{s-1}|\sinh (t)|^{-s}\right)
$$

We now use $\ell_{j} \sinh (t)=\rho$ and $\ell_{j}^{s-1}|\sinh (t)|^{-s}=\ell_{j}^{2 s-1}|\rho|^{-s} \leqslant \varepsilon|\rho|^{s-1}$ if $|\rho|>C \ell_{j}$ with $C$ large enough (depending on $\varepsilon$ ).

Next, consider the case $|\rho|<C \ell_{j}$. We can also write $u_{0}$ and $v_{0}$ under the form (see [Bo, Chapter 5.5])

$$
\begin{gathered}
v_{0}(t)=\frac{\Gamma\left(\frac{1+s}{2}\right)^{2}}{\Gamma\left(\frac{3}{2}\right) \Gamma\left(s-\frac{1}{2}\right)} \sinh (t) F\left(\frac{1+s}{2}, 1-\frac{s}{2}, \frac{3}{2} ;-\sinh (t)^{2}\right), \\
u_{0}(t)=\frac{\Gamma\left(\frac{s}{2}\right)^{2}}{\Gamma\left(\frac{1}{2}\right) \Gamma\left(\frac{1}{2}-s\right)} F\left(\frac{s}{2}, \frac{1-s}{2}, \frac{1}{2} ;-\sinh (t)^{2}\right)
\end{gathered}
$$

and this easily yields the desired estimate.

The proof is complete by noticing that the estimates on $H_{k}, J_{k}$ follow fom this Lemma and (2.41), together with the fact that each $\varphi_{\lambda_{k}}$ is bounded uniformly in $M \backslash \cup_{j} \mathcal{C}_{j}(g)$ for $g \in U_{g_{0}}$ by Lemma 2.5 .

## 3. Gaussian Free Field and Gaussian Multiplicative Chaos

In this section, we shall explain how to give a mathematical sense to the formal measure

$$
\begin{equation*}
F \mapsto \int F(\varphi) e^{-S_{L}(g, \varphi)} D \varphi \tag{3.1}
\end{equation*}
$$

where $S_{L}(g, \varphi)$ is the Liouville functional defined in (2.2), $g$ is a fixed metric on the surface $M$ and $\varphi$ varies among a certain space of functions so that $e^{\gamma \varphi} g$ is parametrizing the conformal class $[g]$ of $g$. This will allow us to define the partition function of Liouville Quantum Field Theory, and in fact $\varphi$ will be a field, i.e. a random function or random distribution, that we will denote by $X_{g}$. The first step is to make sense of the part corresponding to the squared gradient term in $S_{L}(g, \varphi)$, i.e. the formal Gaussian measure

$$
\begin{equation*}
F \mapsto \int F(\varphi) e^{-\frac{1}{4 \pi}\|d \varphi\|_{L^{2}}^{2} D \varphi} \tag{3.2}
\end{equation*}
$$

Classically, we interpret the above field $\varphi$ as a Gaussian Free Field (GFF in short): this is a Gaussian random variable taking values in some space of distributions in the sense of Schwartz. In particular, the field $\varphi$ is not a fairly defined function and giving sense to the term $e^{\gamma \varphi}$ in (2.2) is thus not straightforward, but it can be done through the theory of Gaussian Multiplicative Chaos (GMC in short), which goes back to [Ka].
3.1. Gaussian Free Field. We describe the Gaussian Free Field (GFF) on a compact Riemannian surface $(M, g)$ by using our previous description of the Green function. The definition of the GFF, as well as the definition of its partition function, can be carried out in a direct way ([Du1, She]). Yet, this path is maybe not as pedagogical as following the threads of ideas that have led physicists to our current knowledge of this object, which we try to summarize below in a rather heuristic way to end up with a mathematically sounded definition.

As a warm up, let us quickly recall that the Gaussian measure

$$
(2 \pi)^{-n / 2} \sqrt{\operatorname{det}(A)} e^{-\frac{1}{2}\langle A x, x\rangle} d x
$$

on $\mathbb{R}^{n}$ when $A$ is a positive definite symmetric operator is the law of the random variable $X=\sum_{j=1}^{n} \alpha_{j} \varphi_{j} / \sqrt{\lambda_{j}}$ where $\left(\alpha_{j}\right)_{j}$ are independent Gaussian random variables in $\mathcal{N}(0,1)$ (zero mean and variance 1), and $\left(\varphi_{j}\right)_{j}$ is an orthonormal basis of eigenvectors for $A$ with eigenvalues $\lambda_{j}>0$.

As the GFF is an infinite dimensional Gaussian, it is natural to expect a construction through its projections onto finite dimensional subspaces, on which one can apply the construction described just above. For this, recall that the Laplacian $\Delta_{g}$ has an orthonormal basis of real valued eigenfunctions $\left(\varphi_{j}\right)_{j \in \mathbb{N}_{0}}$ in $L^{2}(M, g)$ with associated eigenvalues $\lambda_{j} \geqslant 0$. We set $\lambda_{0}=0$ and $\varphi_{0}$ to be constant. The Laplacian can thus be seen as a symmetric operator on an infinite dimensional space. Denote $H_{n}$ the finite dimensional space spanned by the first $n$ eigenfunctions $\left(\varphi_{j}\right)_{1 \leqslant j \leqslant n}$ of the Laplacian. Notice that for $\phi=\sum_{j=1}^{n} \widetilde{\alpha}_{j} \varphi_{j}$ we have $\|d \phi\|_{L^{2}}^{2}=\sum_{j=1}^{n} \widetilde{\alpha}_{j}^{2} \lambda_{j}$.

Therefore the projection to $H_{n}$ of the formal measure (3.2) is naturally understood as

$$
\begin{aligned}
\int_{H_{n}} F(\phi) e^{-\frac{1}{4 \pi}\|d \phi\|_{L^{2}}^{2}} D \phi & =\int_{\mathbb{R}^{n}} F\left(\sum_{j=1}^{n} \widetilde{\alpha}_{j} \varphi_{j}\right) \prod_{j=1}^{n}\left(e^{-\frac{1}{4 \pi}\left(\widetilde{\alpha}_{j}\right)^{2} \lambda_{j}} d \widetilde{\alpha}_{j}\right) \\
& =(2 \pi)^{n / 2}\left(\prod_{j=1}^{n} \lambda_{j}\right)^{-1 / 2} \int_{\mathbb{R}^{n}} F\left(\sqrt{2 \pi} \sum_{j=1}^{n} \alpha_{j} \frac{\varphi_{j}}{\sqrt{\lambda_{j}}}\right) \prod_{j=1}^{n}\left(e^{-\frac{\alpha_{j}^{2}}{2}} d \alpha_{j}\right)
\end{aligned}
$$

for appropriate bounded measurable functionals $F$. The total mass of this measure is easily computed and equals $(2 \pi)^{n}\left(\prod_{j=1}^{n} \lambda_{j}\right)^{-1 / 2}$. Furthermore, when renormalized by its total mass, this measure becomes thus a probability measure describing the law of the random function

$$
\begin{equation*}
X_{n}=\sqrt{2 \pi} \sum_{j=1}^{n} \alpha_{j} \frac{\varphi_{j}}{\sqrt{\lambda_{j}}} \tag{3.3}
\end{equation*}
$$

where $\left(\alpha_{j}\right)_{j}$ are independent Gaussian random variables with law $\mathcal{N}(0,1)$.
To obtain the description of the GFF, one has to take the limit $n \rightarrow \infty$. It can be seen [Du1] that the sum (3.3) converges almost surely in the Sobolev space $H^{-s}(M)$ for $s>0$. The total mass $(2 \pi)^{n}\left(\prod_{j=1}^{n} \lambda_{j}\right)^{-1 / 2}$ diverges as $n \rightarrow \infty$ but this is not that much troublesome as it is customary in physics (and mathematics) to remove the diverging terms provided they are "universal enough" (this procedure is called renormalization). Removing the diverging terms should give a limiting total mass equal to $\left(\operatorname{det}^{\prime}\left(\frac{1}{2 \pi} \Delta_{g}\right) / \operatorname{Vol}_{g}(M)\right)^{-1 / 2}$. So far, this is the picture the reader should have in mind to understand the construction of the GFF. Yet, for readers who want to have more details, we stress that renormalizing the product $\prod_{j=1}^{n} \lambda_{j}$ turns out to be very troublesome and slight adaptations are necessary to recover the phenomenology explained above. The reader may consult the paper [BiFe] where these renormalization issues are discussed in further details.

The above formal discussion thus motivates the forthcoming definitions. The Green function $G_{g}\left(x, x^{\prime}\right)$ (with vanishing mean) is a distribution on $M \times M$ which can be written as the series (converging in the sense of distributions)

$$
G_{g}\left(x, x^{\prime}\right)=\sum_{j=1}^{\infty} \frac{\varphi_{j}(x) \varphi_{j}\left(x^{\prime}\right)}{\lambda_{j}}
$$

Let $\left(a_{j}\right)_{j}$ be a sequence of i.i.d. real Gaussians $\mathcal{N}(0,1)$, defined on some probability space $(\Omega, \mathcal{F}, \mathbb{P})$, and define the Gaussian Free Field (GFF) with vanishing mean in the metric $g$ by

$$
\begin{equation*}
X_{g}=\sqrt{2 \pi} \sum_{j \geqslant 1} a_{j} \frac{\varphi_{j}}{\sqrt{\lambda_{j}}} \tag{3.4}
\end{equation*}
$$

as a random variable with values in distributions on $M$, i.e. almost surely $X_{g} \in \mathcal{D}^{\prime}(M)$ (see [Du1] section 4.2 for instance). Notice that for each $\phi \in C^{\infty}(M)$, almost surely, we have $\left\langle X_{g}, \phi\right\rangle=\sqrt{2 \pi} \sum_{j=1}^{\infty} a_{j} \frac{\left\langle\varphi_{j}, \phi\right\rangle}{\sqrt{\lambda_{j}}}$ which is a converging series of random variables as $\mathbb{E}\left(\left\langle X_{g}, \phi\right\rangle^{2}\right)<$ $\infty$. In fact, if $H_{0}^{-s}(M)$ is the kernel of the map $X \mapsto\langle X, 1\rangle_{L^{2}\left(\mathrm{dv}_{g}\right)}$ on the $L^{2}$-based Sobolev space $H^{-s}(M)$ of order $-s \in \mathbb{R}_{-}^{*}$, it is easy to see (see [Du1] again) that $X_{g}$ makes sense as a
random variable with values in $L^{2}\left(\Omega ; H_{0}^{-s}(M)\right)$ for all $s>0$ by using the asymptotic counting function on the eigenvalues $\lambda_{j}$ (i.e. Weyl law). If $\phi_{1}, \phi_{2} \in C^{\infty}(M)$, we have the covariance value

$$
\mathbb{E}\left[\left\langle X_{g}, \phi_{1}\right\rangle \cdot\left\langle X_{g}, \phi_{2}\right\rangle\right]=2 \pi \sum_{j=1}^{\infty} \frac{\left\langle\varphi_{j}, \phi_{1}\right\rangle\left\langle\varphi_{j}, \phi_{2}\right\rangle}{\lambda_{j}}=2 \pi\left\langle G_{g}, \phi_{1} \otimes \phi_{2}\right\rangle .
$$

The covariance is the Green's function when viewed as a distribution: if $\phi_{1} \rightarrow \delta_{x}$ and $\phi_{2} \rightarrow \delta_{x^{\prime}}$ for $x \neq x^{\prime}$ we have $\mathbb{E}\left(\left\langle X_{g}, \phi_{1}\right\rangle \cdot\left\langle X_{g}, \phi_{2}\right\rangle\right) \rightarrow 2 \pi G_{g}\left(x, x^{\prime}\right)$ and we write the formal equality for the covariance of the GFF

$$
\mathbb{E}\left[X_{g}(x) \cdot X_{g}\left(x^{\prime}\right)\right]=2 \pi G_{g}\left(x, x^{\prime}\right)
$$

Notice that the extra $2 \pi$ factor serves to make the field $X_{g}$ have exact logarithmic correlations in view of Lemma 2.1. As in [She, Theorem 2.3], there is a probability mesure $\mathcal{P}$ on $H_{0}^{-s}(M)$ (for some natural $\sigma$-algebra) so that the law of $X_{g}$ is given by $\mathcal{P}$ and for each $\phi \in H^{s}(M)$, $\left\langle X_{g}, \phi\right\rangle$ is a random variable on $\Omega$ with zero mean and variance $2 \pi\left\langle R_{g}(0) \phi, \phi\right\rangle$. The mesure $\mathcal{P}$ represents the Gaussian mesure (3.2) (times the $\sqrt{\operatorname{det}\left(\Delta_{g}\right)}$ term) on the space of functions orthogonal to constants, thus to define (3.2) on the whole $H^{-s}(M)$ space, we shall consider the tensor product $\mathcal{P} \otimes d c$ where $d c$ is the Lebesgue mesure on $\mathbb{R}$ viewed as the 1-dimensional vector space of constant functions on $M$ : in other words, we use the isomorphism

$$
H_{0}^{-s}(M) \times \mathbb{R} \rightarrow H^{-s}(M), \quad(X, c) \mapsto X+c .
$$

to define the mesure $\mathcal{P}^{\prime}$ on $H^{-s}(M)$ as the image of $\mathcal{P} \otimes d c$ by this map. This mesure gives a proper sense to the Gaussian mesure (3.2). We have

Lemma 3.1. The mesure $\mathcal{P}^{\prime}$ on $H^{-s}(M)$ obtained by tensorizing the GFF mesure $\mathcal{P}$ by dc is conformally invariant in the sense that it does not depend on the conformal representative in a conformal class $[g]$.

Proof. Let $\hat{g}=e^{\omega} g$ for some $\omega \in C^{\infty}(M)$. Notice that $H_{0}^{-s}(M)$ depends on $g$, we thus denote it $H_{0}^{-s}(M, g)$ and we denote $\langle\cdot, \cdot\rangle_{g}$ the distribution pairing on $M$ or $M \times M$ induced by the measure $\mathrm{dv}_{g}$. First we claim that the probability law obtained from $\hat{X}_{g}:=X_{g}-c_{\hat{g}}\left(X_{g}\right)$ is the same as that of $X_{\hat{g}}$, if $c_{\hat{g}}\left(X_{g}\right):=\left\langle X_{g}, 1\right\rangle_{\hat{g}} / \operatorname{Vol}_{\hat{g}}(M)=\left\langle X_{g}, e^{\omega}\right\rangle_{g} / \operatorname{Vol}_{\hat{g}}(M)$. The random field $\hat{X}_{g}$ satisfies $\left\langle\hat{X}_{g}, 1\right\rangle_{\hat{g}}=0$ and is thus in the space $H_{0}^{-s}(M, \hat{g})$, moreover $\mathbb{E}\left[\left\langle\hat{X}_{g}, \phi\right\rangle_{\hat{g}}\right]=0$ for all $\phi \in C^{\infty}(M)$. The covariance of $\hat{X}_{g}$ is given by

$$
\begin{aligned}
\mathbb{E}\left[\left\langle\hat{X}_{g}, \phi_{1}\right\rangle_{\hat{g}}\left\langle\hat{X}_{g}, \phi_{2}\right\rangle_{\hat{g}}\right]= & \left\langle G_{g}, \phi_{1} \otimes \phi_{2}\right\rangle_{\hat{g}}+\left(\operatorname{Vol}_{\hat{g}}(M)\right)^{-2}\left\langle G_{g}, 1 \otimes 1\right\rangle_{\hat{g}}\left\langle\phi_{1}, 1\right\rangle_{\hat{g}}\left\langle\phi_{2}, 1\right\rangle_{\hat{g}} \\
& -\left(\operatorname{Vol}_{\hat{g}}(M)\right)^{-1}\left(\left\langle G_{g}, 1 \otimes \phi_{2}\right\rangle_{\hat{g}}\left\langle 1, \phi_{1}\right\rangle_{\hat{g}}+\left\langle G_{g}, \phi_{1} \otimes 1\right\rangle_{\hat{g}}\left\langle 1, \phi_{2}\right\rangle_{\hat{g}}\right) \\
= & \left\langle G_{g}+\alpha 1 \otimes 1-u \otimes 1-1 \otimes u, \phi_{1} \otimes \phi_{2}\right\rangle_{\hat{g}}
\end{aligned}
$$

where $\alpha=\left(\operatorname{Vol}_{\hat{g}}(M)\right)^{-2}\left\langle G_{g}, 1 \otimes 1\right\rangle_{\hat{g}}, u(x)=\int_{M} G_{g}(x, y) \operatorname{dv}_{\hat{g}}(y) / \operatorname{Vol}_{\hat{g}}(M)$. We recognize that this kernel is just the Green function for $\hat{g}$ (this is easily checked) against $\phi_{1} \otimes \phi_{2}$, showing that the correlation of $\hat{X}_{g}$ is that of $X_{\hat{g}}$. Since both random fields are Gaussian, we deduce that the law of $\hat{X}_{g}$ and $X_{\hat{g}}$ are the same and thus for $F \in L^{1}\left(H^{-s}(M), \mathcal{P}^{\prime}\right)$,

$$
\int_{\mathbb{R}} \mathbb{E}\left[F\left(X_{\hat{g}}+c\right)\right] d c=\int_{\mathbb{R}} \mathbb{E}\left[F\left(X_{g}-c_{\hat{g}}\left(X_{g}\right)+c\right)\right] d c=\int_{\mathbb{R}} \mathbb{E}\left[F\left(X_{g}+c\right)\right] d c
$$

by making a change of variables in $c$.

Finally, in view of the discussion above, the measure $F \mapsto \mathbb{E}(F)$ on $H^{-s}(M)$ represents the formal measure

$$
\begin{equation*}
\left.\sqrt{\operatorname{det}^{\prime}\left(\frac{1}{2 \pi} \Delta_{g}\right.}\right) e^{-\frac{1}{4 \pi} \int_{M}|d \varphi|_{g}^{2} \mathrm{dv}_{g}} D \varphi \tag{3.5}
\end{equation*}
$$

and using (2.10), we can write $\sqrt{\operatorname{det}^{\prime}\left(\frac{1}{2 \pi} \Delta_{g}\right)}=(2 \pi)^{\frac{1}{2}\left(1-\frac{\chi(M)}{6}\right)} \sqrt{\operatorname{det}^{\prime}\left(\Delta_{g}\right)}$.
3.2. Gaussian multiplicative chaos. To define quantities like $e^{\gamma X}$ for some $\gamma \in \mathbb{R}$ we will use a renormalization procedure after regularization of the field $X_{g}$. We describe the construction for $g$ hyperbolic and we shall remark that in fact the construction works as well for any conformal metric $\hat{g}=e^{\omega} g$ by using Lemma 2.1.

First, when $\varepsilon>0$ is very small, we define a regularization $X_{g, \varepsilon}$ of $X_{g}$ by averaging on geodesic circles of radius $\varepsilon>0$. Let $x \in M$ and let $\mathcal{C}(x, \varepsilon)$ be the geodesic circle of center $x$ and radius $\varepsilon>0$, and let $\left(f_{x, \varepsilon}^{n}\right)_{n \in \mathbb{N}} \in C^{\infty}(M)$ be a sequence with $\left\|f_{x, \varepsilon}^{n}\right\|_{L^{1}}=1$ which is given by $f_{x, \varepsilon}^{n}=\theta^{n}\left(d_{g}(x, \cdot) / \varepsilon\right)$ where $\theta^{n}(r) \in C_{c}^{\infty}((0,2))$ non-negative supported near $r=1$ such that $f_{x, \varepsilon}^{n} \mathrm{dv}_{g}$ is converging in $\mathcal{D}^{\prime}(M)$ to the uniform probability measure $\mu_{x, \varepsilon}$ on $\mathcal{C}_{g}(x, \varepsilon)$ as $n \rightarrow \infty$ (for $\epsilon$ small enough, the geodesic circles form a $1 d$-manifold and the trace of $g$ along this manifold gives rise to a finite measure, which corresponds to the uniform measure after renormalization so as to have mass 1 , it can also be defined in terms of 1-dimensional Hausdorff measure constructed with the volume form on $M$ and restricted to this geodesic circle). Then we have the standard

Lemma 3.2. The random variable $\left\langle X_{g}, f_{x, \varepsilon}^{n}\right\rangle$ converges to a random variable as $n \rightarrow \infty$, which has a modification $X_{g, \varepsilon}(x)$ with continuous sample paths with respect to $(x, \varepsilon) \in M \times\left(0, \varepsilon_{0}\right)$, with covariance

$$
\mathbb{E}\left[X_{g, \varepsilon}(x) X_{g, \varepsilon}\left(x^{\prime}\right)\right]=2 \pi \int G_{g}\left(y, y^{\prime}\right) d \mu_{x, \varepsilon}(y) d \mu_{x^{\prime}, \varepsilon}\left(y^{\prime}\right)
$$

and we have as $\varepsilon \rightarrow 0$

$$
\begin{equation*}
\mathbb{E}\left[X_{g, \varepsilon}(x)^{2}\right]=-\log (\varepsilon)+W_{g}(x)+o(1) \tag{3.6}
\end{equation*}
$$

where $W_{g}$ is the smooth function on $M$ given by $W_{g}(x)=2 \pi m_{g}(x, x)+\frac{1}{2} \log (2)$ if $m_{g}$ is the smooth function of Lemma 2.1.

Proof. Let us fix $x, \varepsilon$, then if $Y_{n}:=\left\langle X_{g}, f_{x, \varepsilon}^{n}\right\rangle$, it suffices to show that $\mathbb{E}\left(Y_{n} Y_{n^{\prime}}\right)$ has a limit as $\left(n, n^{\prime}\right) \rightarrow \infty$ to prove that $Y_{n}$ is a Cauchy sequence in $L^{2}(\Omega)$. Using Lemma 2.1 (and its notation):

$$
\begin{align*}
\mathbb{E}\left(Y_{n} Y_{n^{\prime}}\right) & =2 \pi \int_{M \times M} G_{g}\left(y, y^{\prime}\right) f_{x, \varepsilon}^{n}(y) f_{x, \varepsilon}^{n^{\prime}}\left(y^{\prime}\right) \mathrm{dv}_{g}(y) \mathrm{dv}_{g}\left(y^{\prime}\right) \\
& =\int_{M \times M}\left(-\log \left(d_{g}\left(y, y^{\prime}\right)\right)+2 \pi m_{g}\left(y, y^{\prime}\right)\right) \theta_{\varepsilon}^{n}\left(d_{g}(x, y)\right) \theta_{\varepsilon}^{\theta^{\prime}}\left(d_{g}\left(x, y^{\prime}\right)\right) \operatorname{dv}_{g}(y) \operatorname{dv}_{g}\left(y^{\prime}\right) . \tag{3.7}
\end{align*}
$$

Clearly the term

$$
\int_{M \times M} m_{g}\left(y, y^{\prime}\right) \theta_{\varepsilon}^{n}\left(d_{g}(x, y)\right) \theta_{\varepsilon}^{n^{\prime}}\left(d_{g}\left(x, y^{\prime}\right)\right) \operatorname{dv}_{g}(y) \operatorname{dv}_{g}\left(y^{\prime}\right)
$$

is uniformly bounded in $\left(n, n^{\prime}, \varepsilon\right)$ and, as $\left(n, n^{\prime}\right) \rightarrow \infty$, it converges, to

$$
\int_{\mathcal{C}(x, \varepsilon)} \int_{\mathcal{C}(x, \varepsilon)} m_{g}\left(y, y^{\prime}\right) d \mu_{x, \varepsilon}(y) d \mu_{x, \varepsilon}\left(y^{\prime}\right)
$$

which in turn is smooth in $x$ and converges, as $\varepsilon \rightarrow 0$, to $m_{g}(x, x)$ uniformly in $x$. For $\varepsilon>0$ small enough, we can use an isometry $\psi$ between a small geodesic ball $B_{g}(x, 3 \varepsilon)$ of radius $3 \varepsilon$ and the ball $B_{\mathbb{H}^{2}}(0,3 \varepsilon)$ in $\mathbb{H}^{2}$ viewed as the disk model, so that the integral (3.7) above reduces to an integral in $B_{g}(x, 3 \varepsilon)$ in both $y, y^{\prime}$. Using the coordinates $z \in \mathbb{H}^{2}$ induced by $\psi$ and (2.20), we get

$$
\begin{gathered}
\int_{M \times M} \log \left(d_{g}\left(y, y^{\prime}\right)\right) \theta_{\varepsilon}^{n}\left(d_{g}(x, y)\right) \theta_{\varepsilon}^{n^{\prime}}\left(d_{g}\left(x, y^{\prime}\right)\right) \mathrm{dv}_{g}(y) \mathrm{dv}_{g}\left(y^{\prime}\right)= \\
\int_{[0,1]^{2} \times[0,2 \pi]^{2}}\left(\log \left|\tanh \left(\frac{r}{2}\right) e^{i\left(\alpha-\alpha^{\prime}\right)}-\tanh \left(\frac{r^{\prime}}{2}\right)\right|+L\right) \theta^{n}\left(\frac{r}{\varepsilon}\right) \theta^{n^{\prime}}\left(\frac{r^{\prime}}{\varepsilon}\right) d \alpha d \alpha^{\prime} \sinh (r) \sinh \left(r^{\prime}\right) d r d r^{\prime} .
\end{gathered}
$$

where $L=L\left(r, r^{\prime}, e^{i \alpha}, e^{i \alpha}\right)$ is continuous and $L(0,0, \cdot, \cdot)=0$. The term involving $L$ is clearly uniformly bounded in $\left(n, n^{\prime}\right)$ and $\varepsilon$ and converges just like for $m_{g}$ above, and its limit as $\varepsilon \rightarrow 0$ is 0 . The part with the $\log$ term is also straightforward to deal with and is also uniformly bounded in ( $n, n^{\prime}$ ) for fixed $\varepsilon>0$ and we get

$$
\begin{gathered}
\int_{[0,1]^{2} \times[0,2 \pi]^{2}} \log \left|\tanh \left(\frac{r}{2}\right) e^{i\left(\alpha-\alpha^{\prime}\right)}-\tanh \left(\frac{r^{\prime}}{2}\right)\right| \theta^{n}\left(\frac{r}{\varepsilon}\right) \theta^{n^{\prime}}\left(\frac{r^{\prime}}{\varepsilon}\right) d \alpha d \alpha^{\prime} \sinh (r) \sinh \left(r^{\prime}\right) d r d r^{\prime} \\
\underset{\left(n, n^{\prime}\right) \rightarrow \infty}{\longrightarrow} \log \left|\tanh \left(\frac{\varepsilon}{2}\right)\right|+\frac{1}{4 \pi^{2}} \int_{[0,2 \pi]^{2}} \log \left|e^{i\left(\alpha-\alpha^{\prime}\right)}-1\right| d \alpha d \alpha^{\prime}=\log \left|\tanh \left(\frac{\varepsilon}{2}\right)\right| .
\end{gathered}
$$

We then have shown the convergence of $\left\langle X_{g}, f_{x, \varepsilon}^{n}\right\rangle$ towards a random variable $\widetilde{X}_{g, \varepsilon}(x)$ in $L^{2}(\Omega)$. To show it has a modification $X_{g, \varepsilon}(x)$ that is sample continuous in $(x, \varepsilon) \in M \times\left(0, \varepsilon_{0}\right)$, it suffices to apply Kolmogorov multi-parameter continuity theorem exactly like in the proof of [DuSh, Prop. 3.1], we do not repeat the argument. The variance $\mathbb{E}\left(X_{g, \varepsilon}(x)^{2}\right)$ is smooth in $x$ and behaves like $-\log (\varepsilon)+\frac{1}{2} \log (2)+2 \pi m_{g}(x, x)+o(1)$ as $\varepsilon \rightarrow 0$, uniformly in $x$.

Next from Lemma 3.2, we will be able to define the Gaussian Multiplicative Chaos (GMC) first considered by Kahane [Ka] in the eighties. The reader may also consult [DuSh, RhVa1, RoVa, Sha] on the topic (in particular, we recommend [Be] for the simplicity of the approach).
Proposition 3.3. 1) Let $\gamma>0$, the random measures $\mathcal{G}_{g, \varepsilon}^{\gamma}:=\varepsilon^{\frac{\gamma^{2}}{2}} e^{\gamma X_{g, \varepsilon}(x)} \operatorname{dv}_{g}(x)$ converge in probability and weakly in the space of Radon measures towards a random measure $\mathcal{G}_{g}^{\gamma}(d x)$. The measure $\mathcal{G}_{g}^{\gamma}(d x)$ is non zero if and only if $\gamma \in(0,2)$.
2) One obtains a non trivial random measure that we will denote $\mathcal{G}_{g}^{2}$ in the case $\gamma=2$ by considering the limit in probability and in the sens of weak convergence of measures of the family of random measures $\mathcal{G}_{g, \varepsilon}^{\gamma}:=(-\ln \varepsilon)^{1 / 2} \varepsilon^{\frac{\gamma^{2}}{2}} e^{\gamma X_{g, \varepsilon}(x)} \mathrm{dv}_{g}(x)$.

Proof. The proof is standard for convolution based regularizations of log-correlated Gaussian fields (first considered in [RoVa], see [Sha] for latest results) in the case $\gamma<2$. The case $\gamma=2$ is treated in [DRV, section 5] in the case of tori, relying on the result proved in [DRSV1, DRSV2] for GFF with Dirichlet boundary conditions. The same argument applies for general compact $2 d$-surfaces up to cosmetic modifications.

Here we give a simple argument in the case $\gamma<\sqrt{2}$ for convenience of readers who are not familiar with GMC. Using the expression (3.6), it suffices to study the convergence of the measures

$$
e^{\gamma X_{g, \varepsilon}(x)-\frac{\gamma^{2}}{2} \mathbb{E}\left[X_{g, \varepsilon}(x)^{2}\right]} e^{\frac{\gamma^{2}}{2} W_{g}(x)} \operatorname{dv}_{g}(x) .
$$

Then by Fubini we directly get for each Borel set $A \subset M$, with $d \sigma=e^{\frac{\gamma^{2}}{2} W_{g}(x)} \operatorname{dv}_{g}(x)$

$$
\mathbb{E}\left[\mathcal{G}_{g, \varepsilon}^{\gamma}(A)\right]=\int_{A} \mathbb{E}\left[e^{\gamma X_{g, \varepsilon}(x)-\frac{\gamma^{2}}{2} \mathbb{E}\left[X_{g, \varepsilon}(x)^{2}\right]}\right] d \sigma(x)=\sigma(A) .
$$

Using that there is $C>1$ such that for all $z \in \mathbb{C},|z|<1$ and $\varepsilon>0$ small

$$
1 / C+|\log (|z|+\varepsilon)| \leqslant \int_{0}^{2 \pi}|\log | z+\varepsilon e^{i \alpha}| | d \alpha \leqslant C+|\log (|z|+\varepsilon)|
$$

then the arguments in the proof of Lemma 3.2 and the expression (2.20) imply that there is $C^{\prime}$ such that

$$
\begin{equation*}
1 / C^{\prime}+\left|\log \left(d_{g}\left(x^{\prime}, x\right)+\varepsilon\right)\right| \leqslant \mathbb{E}\left[X_{g, \varepsilon}(x) X_{g, \varepsilon^{\prime}}\left(x^{\prime}\right)\right] \leqslant C^{\prime}+\left|\log \left(d_{g}\left(x^{\prime}, x\right)+\varepsilon\right)\right| . \tag{3.8}
\end{equation*}
$$

for all $\varepsilon^{\prime} \leqslant \varepsilon$, and all $x, x^{\prime} \in M$. In particular we get by using Fubini and the fact that $X_{g}$ is a Gaussian free field

$$
\begin{aligned}
\mathbb{E}\left[\mathcal{G}_{g, \varepsilon}^{\gamma}(A)^{2}\right] & =\mathbb{E}\left[\left(\int_{A} e^{\gamma X_{g, \varepsilon}(x)-\frac{\gamma^{2}}{2} \mathbb{E}\left[X_{g, \varepsilon}(x)^{2}\right]} d \sigma(x)\right)^{2}\right] \\
& =\mathbb{E}\left[\int_{A} \int_{A} e^{\gamma\left(X_{g, \varepsilon}(x)+X_{g, \varepsilon}\left(x^{\prime}\right)\right)-\frac{\gamma^{2}}{2}\left(\mathbb{E}\left[X_{g, \varepsilon}(x)^{2}\right]+\mathbb{E}\left[X_{g, \varepsilon}\left(x^{\prime}\right)^{2}\right]\right)} d \sigma(x) d \sigma\left(x^{\prime}\right)\right] \\
& =\int_{A} \int_{A} e^{\gamma^{2} \mathbb{E}\left[X_{g, \varepsilon}(x) X_{g, \varepsilon}\left(x^{\prime}\right)\right]} d \sigma(x) d \sigma\left(x^{\prime}\right)
\end{aligned}
$$

which converges to $\int_{A} \int_{A} e^{\gamma^{2} 2 \pi G_{g}\left(x, x^{\prime}\right)} d \sigma(x) d \sigma\left(x^{\prime}\right)<\infty$ as $\varepsilon \rightarrow 0$ by using (3.8) and Lebesgue theorem - the condition $\gamma^{2}<2$ appear here due to the $\log$ divergence of $2 \pi G_{g}\left(x, x^{\prime}\right)$ at $x=x^{\prime}$, see Lemma 2.1. A similar argument and (3.8) also show that $\mathbb{E}\left[\left(\mathcal{G}_{g, \varepsilon}^{\gamma}(A)-\mathcal{G}_{g, \varepsilon^{\prime}}^{\gamma}(A)\right)^{2}\right] \rightarrow 0$ if $\left(\varepsilon, \varepsilon^{\prime}\right) \rightarrow 0$, thus $\mathcal{G}_{g, \varepsilon}^{\gamma}(A)$ is a Cauchy sequence, which therefore converges in $L^{2}(\Omega)$ to a random variable $Z(A)$, of mean $\sigma(A)$. By standard arguments, $\mathcal{G}_{g, \varepsilon}^{\gamma}(d x)$ converges to a random measure $\mathcal{G}_{g}^{\gamma}$ satisfying $\mathbb{E}\left[\mathcal{G}_{g}^{\gamma}(A)\right]=\sigma(A)$. The case $\gamma \in[\sqrt{2}, 2)$ is more complicated and several methods have been proposed in the literature. We refer to [Be] for a simple argument.

In fact, the whole construction above is not so particular to choosing the hyperbolic metric: indeed it uses only the fact that the covariance of $X_{g}$ is the Green function, the fact that near the diagonal $2 \pi G_{g}\left(x, x^{\prime}\right)=-\log d_{g}\left(x, x^{\prime}\right)+F\left(x, x^{\prime}\right)$ with $F$ continuous, and finally the fact that in local isothermal coordinates $z$ so that $g=e^{2 f(z)}|d z|^{2}$

$$
\log d_{g}\left(z, z^{\prime}\right)=\log \left|z-z^{\prime}\right|+\mathcal{O}(1) .
$$

This allows to define a random measure $\mathcal{G}_{\hat{g}}^{\gamma}$ just as above for any other metric $\hat{g}=e^{\omega} g$ conformal to the hyperbolic metric $g$. For later purpose we will need to make the following observation. If $\hat{g}=e^{\omega} g$, define

$$
\begin{equation*}
\hat{X}_{g, \varepsilon}(x):=\lim _{\varepsilon \rightarrow 0}\left\langle X_{g}, \hat{f}_{x, \varepsilon}^{n}\right\rangle_{\hat{g}} \tag{3.9}
\end{equation*}
$$

for each $x \in M$ where $\hat{f}_{x, \varepsilon}^{n}:=\theta^{n}\left(d_{\hat{g}}(x, \cdot) / \varepsilon\right)$ with $\theta^{n}$ like above, so that $\hat{f}_{x, \varepsilon}^{n} \operatorname{dv}_{\hat{g}}$ converge as $n \rightarrow \infty$ to the uniform probability measure $\hat{\mu}_{x, \varepsilon}$ on the geodesic circle $\mathcal{C}_{\hat{g}}(x, \varepsilon)$ of center $x$ and
radius $\varepsilon$ with respect to $\hat{g}$. In isothermal coordinates at $x$ so that $z=0$ correspond to the point $x$ and the metric is $g=|d z|^{2} / \operatorname{Im}(z)^{2}$, the circle $\mathcal{C}_{\hat{g}}(x, \varepsilon)$ is parametrized by

$$
\varepsilon e^{-\frac{1}{2} \omega(z)+\varepsilon h_{\varepsilon}(\alpha)} e^{i \alpha}, \alpha \in[0,2 \pi]
$$

for some continuous function $h_{\varepsilon}(\alpha)$ uniformly bounded in $\varepsilon$. Then one has

$$
\mathbb{E}\left(\hat{X}_{g, \varepsilon}(x) \hat{X}_{g, \varepsilon}\left(x^{\prime}\right)\right)=2 \pi \int G_{g}\left(y, y^{\prime}\right) d \hat{\mu}_{x, \varepsilon}(y) d \hat{\mu}_{x^{\prime}, \varepsilon}\left(y^{\prime}\right)
$$

and by the arguments in the proof of Lemma 3.2, we have as $\varepsilon \rightarrow 0$

$$
\begin{equation*}
\mathbb{E}\left(\hat{X}_{g, \varepsilon}(x)^{2}\right)=-\log (\varepsilon)+W_{g}(x)+\frac{1}{2} \omega(x)+o(1) . \tag{3.10}
\end{equation*}
$$

Then by the arguments of Proposition 3.3, the random measure (add an extra push $(-\ln \varepsilon)^{1 / 2}$ when $\gamma=2$ )

$$
\begin{equation*}
\hat{\mathcal{G}}_{g, \varepsilon}^{\gamma}:=\varepsilon^{\frac{\gamma^{2}}{2}} e^{\gamma \hat{X}_{g, \varepsilon}(x)} \mathrm{dv}_{\hat{g}}(x) \tag{3.11}
\end{equation*}
$$

converges weakly as $\varepsilon \rightarrow 0$ to some measure $\hat{\mathcal{G}}_{g}^{\gamma}$ which satisfies

$$
\begin{equation*}
d \hat{\mathcal{G}}_{g}^{\gamma}(x)=e^{\left(1+\frac{\gamma^{2}}{4}\right) \omega(x)} d \mathcal{G}_{g}^{\gamma}(x) . \tag{3.12}
\end{equation*}
$$

## 4. Liouville Quantum field theory with fixed modulus

In this section we define Liouville Quantum Field Theory (LQFT) with fixed conformal class (also called modulus) and describe its main properties. It follows the approach of [DKRV] in the case of the Riemann sphere. Liouville Quantum Gravity (LQG) with fixed genus is a sum ("partition function") over all possible metrics on a surface with fixed genus. The space of metrics splits into conformal classes and we want to decompose the partition function accordingly. Each conformal class has a unique hyperbolic metric, which plays the role of a background metric.
4.1. Axiomatic of CFT. Here we give a brief account of the axiomatic of Conformal Field Theories in order to motivate the forthcoming results. Our purpose will then be to construct the quantum Liouville theory and show that it satisfies this axiomatic. The reader is referred to [Ga] for more details related to this formalism.

A CFT (on the surface $(M, g)$ ) is described by its partition function $Z(g)$ as well as the correlation functions of the (spinless) primary fields $\left(\theta_{i}\right)_{i \in I}$ denoted by

$$
Z\left(g, \theta_{i_{1}}\left(x_{1}\right), \ldots, \theta_{i_{n}}\left(x_{n}\right)\right)
$$

where $n \geqslant 1,\left\{i_{1}, \ldots, i_{n}\right\} \in I$ and $x_{1}, \ldots, x_{n}$ are arbitrary points on $M$. Let us just roughly say that a CFT is supposed to give sense to "random fields" defined on $M$, here the primary fields $\left(\theta_{i}\right)_{i}$, and the correlation functions can be thought of as the cumulants of these random fields. These correlation functions are supposed to satisfy the following conditions:

- Diffeomorphism covariance: for any orientation preserving diffeomorphism $\psi$

$$
\begin{align*}
Z(g) & =Z\left(\psi^{*} g\right)  \tag{4.1}\\
Z\left(g, \theta_{i_{1}}\left(\psi\left(x_{1}\right)\right), \ldots, \theta_{i_{n}}\left(\psi\left(x_{n}\right)\right)\right) & =Z\left(\psi^{*} g, \theta_{i_{1}}\left(x_{1}\right), \ldots, \theta_{i_{n}}\left(x_{n}\right)\right) \tag{4.2}
\end{align*}
$$

- Conformal anomaly: for any smooth function $\varphi$ on $M$

$$
\begin{align*}
\ln \frac{Z\left(e^{\varphi} g\right)}{Z(g)} & =\frac{\mathbf{c}}{96 \pi} \int_{M}\left(\left|d_{g} \varphi\right|_{g}^{2}+2 K_{g} \varphi\right) \mathrm{dvol}_{g}  \tag{4.3}\\
\ln \frac{Z\left(e^{\varphi} g, \theta_{i_{1}}\left(x_{1}\right), \ldots, \theta_{i_{n}}\left(x_{n}\right)\right)}{Z\left(g, \theta_{i_{1}}\left(x_{1}\right), \ldots, \theta_{i_{n}}\left(x_{n}\right)\right)} & =-\sum_{k=1}^{n} \Delta_{i_{k}} \varphi\left(x_{k}\right)+\frac{\mathbf{c}}{96 \pi} \int_{M}\left(\left|d_{g} \varphi\right|_{g}^{2}+2 K_{g} \varphi\right) \mathrm{dvol}_{g} \tag{4.4}
\end{align*}
$$

where the constant $\mathbf{c}$ is the so-called central charge of the CFT and the scalars $\left(\Delta_{i}\right)_{i \in I}$ are called the conformal weights of the primary fields.

One of the interesting feature of CFTs is their strong algebraic structure, which make them fall under the scope of techniques for integrable systems, leading to the possibility of obtaining exact expressions for the correlation functions, which is the main goal of a Quantum Field Theory.
4.2. The partition function of LQFT. The first step is to describe LQFT with fixed modulus. LQFT will describe the probability law of some random conformal factor, that is we consider the random metrics $e^{\gamma X} g$ if $g$ is a fixed metric and $X$ is a random function. The law of $X$ will be mathematically described by the Feynmann type measure (3.1). So, let $g \in \operatorname{Met}(M)$ be a fixed metric on $M$. The mathematical definition of the measure of LQFT (i.e. (3.1)) is the following. For $F: H^{-s}(M) \rightarrow \mathbb{R}$ (with $s>0$ ) a bounded continuous functional, set for $\gamma \in(0,2]$

$$
\begin{align*}
\Pi_{\gamma, \mu}(g, F):= & \left(\operatorname{det}^{\prime}\left(\Delta_{g}\right) / \operatorname{Vol}_{g}(M)\right)^{-1 / 2}  \tag{4.5}\\
& \times \int_{\mathbb{R}} \mathbb{E}\left[F\left(c+X_{g}\right) \exp \left(-\frac{Q}{4 \pi} \int_{M} K_{g}\left(c+X_{g}\right) \mathrm{dv}_{g}-\mu e^{\gamma c} \mathcal{G}_{g}^{\gamma}(M)\right)\right] d c .
\end{align*}
$$

This quantity, if it is finite, gives a mathematical sense to the formal integral

$$
\int F(\varphi) e^{-S_{L}(g, \varphi)} D \varphi
$$

where $S_{L}(g, \varphi)$ is the Liouville action (2.2). The partition function is the total mass of this measure, i.e $\Pi_{\gamma, \mu}(\hat{g}, 1)$.

Proposition 4.1. For $g \in \operatorname{Met}(M)$ and $\gamma \in(0,2]$, we have $0<\Pi_{\gamma, \mu}(g, 1)<+\infty$ and the mapping

$$
F \in C_{b}\left(H^{-s}(M), \mathbb{R}\right) \mapsto \Pi_{\gamma, \mu}(g, F)
$$

defines a positive finite measure. When renormalized by its total mass, it describes the law of a random variable living in $H^{-s}(M)$ called the Liouville field. When $g \in \operatorname{Met}_{-1}(M)$ is hyperbolic, we further have

$$
\begin{equation*}
\Pi_{\gamma, \mu}(g, 1)=\left(\frac{\operatorname{det}^{\prime}\left(\Delta_{g}\right)}{\operatorname{Vol}_{g}(M)}\right)^{-1 / 2} \gamma^{-1} \mu^{\frac{Q \chi(M)}{\gamma}} \Gamma\left(-\frac{Q \chi(M)}{\gamma}\right) \mathbb{E}\left[\mathcal{G}_{g}^{\gamma}(M)^{\frac{Q \chi(M)}{\gamma}}\right] \tag{4.6}
\end{equation*}
$$

where $\Gamma(z)$ is the standard Euler Gamma function.
Proof. The proof of this proposition follows the same lines as in [DKRV, section 3.1]. We consider the case of a metric $g \in \operatorname{Met}_{-1}(M)$, since the general case follows from that case, as
is explained below in Proposition 4.4 for the correlations functions. In constant curvature, the Gauss-Bonnet theorem entails

$$
\frac{Q}{4 \pi} \int_{M} K_{g}\left(c+X_{g}\right) \mathrm{dv}_{g}=Q c \chi(M)
$$

where $\chi(M)$ is the Euler characteristic of $M$ and we get

$$
\Pi_{\gamma, \mu}(g, F)=\left(\frac{\operatorname{det}^{\prime}\left(\Delta_{g}\right)}{\operatorname{Vol}_{g}(M)}\right)^{-1 / 2} \int_{\mathbb{R}} e^{-Q c \chi(M)} \mathbb{E}\left[\exp \left(-\mu e^{\gamma c} \mathcal{G}_{g}^{\gamma}(M)\right)\right] d c
$$

After inverting expectation and integration, and using the change of variables $y=\mu e^{\gamma c} \mathcal{G}_{\hat{g}}^{\gamma}(M)$, we get (4.6). Finiteness of this quantity is ensured by the fact that GMC has finite moments of negative orders as $\chi(M)<0$ - finiteness of negative moments is proved for example in [RoVa, Proposition 3.6] for $\gamma<2$ and in [DRSV2, ] in the case $\gamma=2$.
4.2.1. Conformal anomaly and diffeomorphism invariance. Here we investigate the symmetries of the measure (4.5) and in particular how the partition function reacts to changes of background metrics. The following proposition is the quantum counterpart of (2.3).

Proposition 4.2. (Conformal anomaly) Let $Q=\frac{\gamma}{2}+\frac{2}{\gamma}$ with $\gamma \in(0,2]$ and $g \in \operatorname{Met}_{-1}(M)$ be a hyperbolic metric on $M$. The partition function satisfies the following conformal anomaly: if $\hat{g}=e^{\omega} g$ for some $\omega \in C^{\infty}(M)$, we have

$$
\Pi_{\gamma, \mu}(\hat{g}, F)=\Pi_{\gamma, \mu}\left(g, F\left(\cdot-\frac{Q}{2} \omega\right)\right) \exp \left(\frac{1+6 Q^{2}}{96 \pi} \int_{M}\left(|d \omega|_{g}^{2}+2 K_{g} \omega\right) \mathrm{dv}_{g}\right) .
$$

Proof. We focus on the integral part in (4.5) (and hence let the determinant of Laplacian apart as its contribution is clear from (2.10)). First, by Lemma 3.1, we can replace $X_{\hat{g}}$ by $X_{g}$ in the expression defining $\Pi_{\gamma, \mu}(\hat{g}, F)$ and are thus left with considering the following quantity (with $\hat{\mathcal{G}}_{g}^{\gamma}$ is the measure defined by (3.11))

$$
A_{\varepsilon}:=\int_{\mathbb{R}} \mathbb{E}\left[F\left(c+X_{g}\right) \exp \left(-\frac{Q}{4 \pi} \int_{M} K_{\hat{g}}\left(c+X_{g}\right) \mathrm{dv}_{\hat{g}}-\mu e^{\gamma c} \hat{\mathcal{G}}_{g}^{\gamma}(M)\right)\right] d c .
$$

By (2.1), the term $-\frac{Q c}{4 \pi} \int_{M} K_{\hat{g}} \mathrm{dv}_{\hat{g}}$ can be written as $-Q c \chi(M)$ where $\chi(M)$ is the Euler characteristic. Define the Gaussian random variable

$$
Y:=-\frac{Q}{4 \pi} \int_{M} K_{\hat{g}} X_{g} \operatorname{dv}_{\hat{g}}=-\frac{Q}{4 \pi}\left\langle X_{g}, K_{\hat{g}} e^{\omega}\right\rangle_{g} .
$$

Let $R_{g}(0)$ be the resolvent operator whose Schwartz kernel is $G_{g}$ with respect to $\mathrm{dv}_{g}$. Since $K_{\hat{g}} e^{\omega}=\Delta_{g} \omega+K_{g}$, we compute, using that $R_{g}(0) K_{g}=0\left(\right.$ as $\left.K_{g}=-2\right)$,

$$
\begin{aligned}
\mathbb{E}\left[\left\langle X_{g}, K_{\hat{g}} e^{\omega}\right\rangle_{g}^{2}\right] & =2 \pi\left\langle G_{g},\left(\Delta_{g} \omega+K_{g}\right) \otimes\left(\Delta_{g} \omega+K_{g}\right)\right\rangle_{g} \\
& =2 \pi\left\langle\omega-\frac{\langle\omega, 1\rangle_{g}}{\operatorname{Vol}_{g}(M)}, \Delta_{g} \omega-2\right\rangle_{g}=2 \pi \int_{M}|d \omega|_{g}^{2} \mathrm{dv}_{g}
\end{aligned}
$$

and similarly we have

$$
\mathbb{E}\left[Y X_{g}\right]=-\frac{Q}{2} R_{g}(0)\left(K_{\hat{g}} e^{\omega}\right)=-\frac{Q}{2}\left(\omega-c_{g}(\omega)\right) .
$$

if $c_{g}(\omega):=\frac{\langle\omega, 1\rangle_{g}}{\operatorname{Vol}_{g}(M)}$. Thus we get

$$
\begin{equation*}
\frac{1}{2} \mathbb{E}\left[Y^{2}\right]=\frac{Q^{2}}{16 \pi} \int_{M}|d \omega|_{g}^{2} \mathrm{dv}_{g}, \quad \mathbb{E}\left[Y X_{g}\right]=-\frac{Q}{2}\left(\omega-c_{g}(\omega)\right) . \tag{4.7}
\end{equation*}
$$

Therefore by applying Girsanov transform to the random variable $Y$, we can rewrite

$$
A=\int_{\mathbb{R}} e^{\frac{1}{2} \mathbb{E}\left[Y^{2}\right]-Q c \chi(M)} \mathbb{E}\left[F\left(c+X_{g}+\mathbb{E}\left(Y X_{g}\right)\right) \exp \left(-\mu e^{\gamma\left(c+\frac{Q}{2} c_{g}(\omega)\right)} \int_{M} e^{-\frac{\gamma Q}{2} \omega} d \hat{\mathcal{G}}_{g}^{\gamma}\right)\right] d c .
$$

With the help of the relation (3.12) and $Q=\frac{\gamma}{2}+\frac{2}{\gamma}$, we see that $\int_{M} e^{-\frac{\gamma Q}{2} \omega} d \hat{\mathcal{G}}_{g}^{\gamma}=\mathcal{G}_{g}^{\gamma}(M)$. Using (4.7), $A$ can be written as

$$
\left.A=\int_{\mathbb{R}} e^{\frac{Q^{2}}{16 \pi}\|d \omega\|_{L_{g}^{2}}^{2}-Q c \chi(M)} \mathbb{E}\left[F\left(c+X_{g}-\frac{Q}{2} \omega+\frac{Q}{2} c_{g}(\omega)\right)\right) \exp \left(-\mu e^{\gamma\left(c+\frac{Q}{2} c_{g}(\omega)\right)} \mathcal{G}_{g}^{\gamma}(M)\right)\right] d c
$$

It remains to make the change of variable $c \rightarrow c-\frac{Q}{2} c_{g}(\omega)$ and we deduce that

$$
A=\int_{\mathbb{R}} e^{\frac{Q^{2}}{16 \pi}\|d \omega\|_{L_{g}^{2}}^{2}-Q c \chi(M)+\frac{1}{2} Q^{2} \chi(M) c_{g}(\omega)} \mathbb{E}\left[F\left(c+X_{g}-\frac{Q}{2} \omega\right) \exp \left(-\mu e^{\gamma c} \mathcal{G}_{g}^{\gamma}(M)\right)\right] d c .
$$

Since $K_{g}=-2$ and $\operatorname{Vol}_{g}(M)=-2 \pi \chi(M)$ we have

$$
-\frac{Q}{4 \pi} \int_{M} K_{g}\left(c+X_{g}\right) \mathrm{dv}_{g}=-Q c \chi(M), \quad c_{g}(\omega) \chi(M)=\frac{1}{4 \pi} \int_{M} K_{g} \omega \operatorname{dv}_{g}
$$

which shows that $A=\Pi_{\gamma, \mu}\left(g, F\left(\cdot-\frac{Q}{2} \omega\right)\right) \sqrt{\operatorname{det}^{\prime}\left(\Delta_{g}\right) / \operatorname{Vol}_{g}(M)} e^{\frac{6 Q^{2}}{} e^{96 \pi} \int_{M}\left(|d \omega|_{g}^{2}+2 K_{g} \omega\right) \mathrm{dv}{ }_{g}}$. Combining with (2.10), the proof is complete.

The constant $c_{L}:=1+6 Q^{2}$ describing the conformal anomaly is called the central charge of the Liouville Theory. Since all the objects in the construction of the Gaussian Free Field and the Gaussian multiplicative chaos are geometric (defined in a natural way from the metric), it is direct to get the following Diffeomorphism invariance:

Proposition 4.3. (Diffeomorphism invariance) Let $g \in \operatorname{Met}(M)$ be a metric on $M$ and let $\psi: M \rightarrow M$ be an orientation preserving diffeomorphism. Then we have for each bounded measurable $F: H^{-s}(M) \rightarrow \mathbb{R}$ with $s>0$

$$
\Pi_{\gamma, \mu}\left(\psi^{*} g, F\right)=\Pi_{\gamma, \mu}(g, F(\cdot \circ \psi)) .
$$

Proof. This follows directly from the fact that all the object considered in the construction of the measure are natural with respect to the metric $g$, thus invariant by isometries: more precisely, it follows from the identities

$$
G_{\psi^{*} g}(x, y)=G_{g}(\psi(x), \psi(y)), \quad K_{\psi^{*} g}(x)=K_{g}(\psi(x)), \quad X_{\psi^{*} g} \stackrel{l a w}{=} X_{g} \circ \psi .
$$

which are obvious.
The two above results show that the axioms (4.1)+(4.3) are satisfied with central charge $\mathbf{c}_{L}=1+6 Q^{2}$. Yet we still have to define the primary fields and their correlation functions. This is the purpose of the next subsection.
4.3. The correlation functions. The correlation functions of LQFT can be thought of as the exponential moments $e^{\alpha \varphi(x)}$ of the random function $\varphi$, the law of which is ruled by the path integral (3.1), evaluated at some location $x \in M$ with weight $\alpha$. Yet, the field $\varphi$ is not a fairly defined function (it belongs to $H^{-s}(M)$ ) so that the construction requires some care.

As before let $g \in \operatorname{Met}(M)$. We fix $n$ points $x_{1}, \ldots, x_{n}(n \geqslant 0)$ on $M$ with respective associated weights $\alpha_{1}, \ldots, \alpha_{n} \in \mathbb{R}$. We denote $\mathbf{x}=\left(x_{1}, \ldots, x_{n}\right)$ and $\boldsymbol{\alpha}=\left(\alpha_{1}, \ldots, \alpha_{n}\right)$. The rigorous definition of the primary fields will require a regularization scheme. We introduce the following $\varepsilon$-regularized functional

$$
\begin{align*}
& \Pi_{\gamma, \mu}^{\mathrm{x}, \boldsymbol{\alpha}}(g, F, \varepsilon):=\left(\operatorname{det}^{\prime}\left(\Delta_{g}\right) / \operatorname{Vol}_{g}(M)\right)^{-1 / 2}  \tag{4.8}\\
& \int_{\mathbb{R}} \mathbb{E}\left[F\left(c+X_{g}\right)\left(\prod_{i} V_{g, \varepsilon}^{\alpha_{i}}\left(x_{i}\right)\right) \exp \left(-\frac{Q}{4 \pi} \int_{M} K_{g}\left(c+X_{g}\right) \operatorname{dv}_{g}-\mu e^{\gamma c} \mathcal{G}_{g, \varepsilon}^{\gamma}(M)\right)\right] d c
\end{align*}
$$

where we have set, for fixed $\alpha \in \mathbb{R}$ and $x \in M$,

$$
V_{g, \varepsilon}^{\alpha}(x)=\varepsilon^{\alpha^{2} / 2} e^{\alpha\left(c+X_{g, \epsilon}(x)\right)} .
$$

Here the regularization is the one described in Lemma 3.2. Such quantities are called vertex operators. Notice that $V_{g, \epsilon}^{\alpha}$ also depends on the variable $c$ but we have dropped this dependence in the notations.

Then, the point is to determine whether the limit

$$
\Pi_{\gamma, \mu}^{\mathbf{x}, \boldsymbol{\alpha}}(g, F):=\lim _{\varepsilon \rightarrow 0} \Pi_{\gamma, \mu}^{\mathbf{x}, \boldsymbol{\alpha}}(g, F, \varepsilon)
$$

exists and defines a non trivial functional on those mappings $F: H^{-s}(M) \rightarrow \mathbb{R}$. If it does, the quantity $\Pi_{\gamma, \mu}^{\mathrm{x}, \boldsymbol{\alpha}}(g):=\Pi_{\gamma, \mu}^{\mathrm{X}, \boldsymbol{\alpha}}(g, 1)$ stands for the $n$-point correlation function of the primary fields $\left(e^{\alpha_{i} \varphi}\right)_{1 \leqslant i \leqslant n}$ respectively evaluated at $\left(x_{i}\right)_{1 \leqslant i \leqslant n}$. Furthermore, another quantity of interest is the probability law on $H^{-s}(M)$ defined by the measure

$$
F \in C_{b}\left(H^{-s}(M)\right) \mapsto \Pi_{\gamma, \mu}^{\mathbf{x}, \alpha}(g, F) / \Pi_{\gamma, \mu}^{\mathrm{X}, \alpha}(g),
$$

which describes the law of some formal "random function" (it is in fact a distribution).
We obtain a result similar to [DKRV] (done for the sphere).
Proposition 4.4. Let $\mathbf{x}=\left(x_{1}, \ldots, x_{n}\right) \in M^{n}$ and $\alpha=\left(\alpha_{1}, \ldots, \alpha_{n}\right) \in \mathbb{R}^{n}$. Then for all bounded continuous functionals $F: h \in H^{-s}(M) \rightarrow F(h) \in \mathbb{R}$ with $s>0$, the limit

$$
\Pi_{\gamma, \mu}^{\mathbf{x}, \alpha}(g, F):=\lim _{\varepsilon \rightarrow 0} \Pi_{\gamma, \mu}^{\mathbf{x}, \alpha}(g, F, \varepsilon),
$$

exists and is finite with $\Pi_{\gamma, \mu}^{\mathbf{X}, \alpha}(g, 1)>0$, if and only if:

$$
\begin{align*}
& \sum_{i} \alpha_{i}+2 Q(\mathbf{g}-1)>0,  \tag{4.9}\\
& \forall i, \quad \alpha_{i}<Q \tag{4.10}
\end{align*}
$$

In the case $g \in \operatorname{Met}_{-1}(M)$, we have the following expression

$$
\Pi_{\gamma, \mu}^{\mathbf{x}, \alpha}(g)=\frac{e^{C(\mathbf{x})} \mu^{-\sum_{i} \alpha_{i}-2 Q(\mathbf{g}-1)} \Gamma\left(\sum_{i} \alpha_{i}+2 Q(\mathbf{g}-1)\right) \mathbb{E}\left[\mathcal{G}_{g, \mathbf{x}, \alpha}^{\gamma}(M)^{-\frac{\sum_{i} \alpha_{i}+2 Q(\mathbf{g}-1)}{\gamma}}\right]}{\sqrt{\operatorname{det}^{\prime}\left(\Delta_{g}\right) / \operatorname{Vol}_{g}(M)}}
$$

where $\Gamma$ is Euler gamma function and, if $W_{g}$ is the function appearing in Lemma 3.2,

$$
\begin{gather*}
\mathcal{G}_{g, \mathbf{x}, \alpha}^{\gamma}(d x):=e^{\gamma \sum_{i} \alpha_{i} 2 \pi G_{g}\left(x_{i}, x\right)} \mathcal{G}_{g}^{\gamma}(d x), \\
C(\mathbf{x}):=\sum_{i} \frac{\alpha_{i}^{2}}{2} W_{g}\left(x_{i}\right)+2 \pi \sum_{i<j} \alpha_{i} \alpha_{j} G_{g}\left(x_{i}, x_{j}\right) . \tag{4.11}
\end{gather*}
$$

Proof. The argument goes essentially as in the proof of [DKRV, Theorems $3.2 \& 3.4$ ], while having in mind that the Gauss-Bonnet theorem is (2.1) on general compact Riemann surfaces. We recall the main steps. It suffices to prove the claim for $F=1$. By Lemma 3.1, we can replace $X_{\hat{g}}$ by $X_{g}$ in the expression defining (4.8) - in particular $V_{\hat{g}, \varepsilon}^{\alpha_{i}}\left(x_{i}\right)$ becomes $\hat{V}_{g, \varepsilon}^{\alpha_{i}}\left(x_{i}\right):=$ $\varepsilon^{\alpha_{i}^{2} / 2} e^{\alpha_{i}\left(c+\hat{X}_{g, \epsilon}\left(x_{i}\right)\right)}$. First we notice by (3.10) that

$$
\hat{V}_{g, \varepsilon}^{\alpha_{i}}\left(x_{i}\right)=e^{\alpha_{i} c+\frac{\alpha_{i}^{2}}{4} \varphi\left(x_{i}\right)+\frac{\alpha_{i}^{2}}{2} W_{g}\left(x_{i}\right)} e^{\alpha_{i} \hat{X}_{g, \varepsilon}\left(x_{i}\right)-\frac{\alpha_{i}^{2}}{2} \mathbb{E}\left[\hat{X}_{g, \varepsilon}\left(x_{i}\right)^{2}\right]}(1+o(1))
$$

as $\varepsilon \rightarrow 0$, with the remainder being deterministic. Here we have used the notation $\hat{X}_{g, \varepsilon}\left(x_{i}\right)=$ $\left\langle X_{g}, \hat{\mu}_{x_{i}, \varepsilon}\right\rangle$ as before, where $\hat{\mu}_{x_{i}, \varepsilon}$ is the uniform probability measure on the Riemannian circle $\mathcal{C}_{\hat{g}}\left(x_{i}, \varepsilon\right)$. Then applying Girsanov transform to

$$
A_{\varepsilon}:=\int_{\mathbb{R}} \mathbb{E}\left[\left(\prod_{i} \hat{V}_{g, \varepsilon}^{\alpha_{i}}\left(x_{i}\right)\right) \exp \left(-\frac{Q}{4 \pi} \int_{M} K_{\hat{g}}\left(c+X_{g}\right) \operatorname{dv}_{\hat{g}}-\mu e^{\gamma^{\gamma}} \hat{\mathcal{G}}_{g, \varepsilon}^{\gamma}(M)\right)\right] d c
$$

we get

$$
A_{\varepsilon}=e^{C_{\varepsilon}(\mathbf{x})} \int_{\mathbb{R}} e^{c\left(\sum_{i} \alpha_{i}-Q \chi(M)\right)} \mathbb{E}\left[\exp \left(-\frac{Q}{4 \pi}\left\langle X_{g}, K_{\hat{g}}\right\rangle_{\hat{g}}-\mu e^{\gamma c} Z_{\varepsilon}\right)\right] d c(1+o(1))
$$

where

$$
\begin{gathered}
Z_{\varepsilon}:=\varepsilon^{\frac{\gamma^{2}}{2}} \int_{M} e^{\gamma\left(\hat{X}_{g, \varepsilon}+H_{g, \varepsilon}\right)} \mathrm{dv}_{\hat{g}} \\
H_{g, \varepsilon}(x):=\sum_{i} 2 \pi \alpha_{i} \int_{\mathcal{C}_{\hat{g}}\left(x_{i}, \varepsilon\right)} G_{g}(y, x) d \hat{\mu}_{x_{i}, \varepsilon}(y), \\
C_{\varepsilon}(\mathbf{x}):=2 \pi \sum_{i \neq j} \alpha_{i} \alpha_{j} G_{g}\left(x_{i}, x_{j}\right)-\frac{Q}{4 \pi} \int_{M} K_{\hat{g}} H_{g, \varepsilon} \operatorname{dv}_{\hat{g}}+\sum_{i} \frac{\alpha_{i}^{2}}{4}\left(\varphi\left(x_{i}\right)+2 W_{g}\left(x_{i}\right)\right) .
\end{gathered}
$$

Notice that, since $K_{\hat{g}} \operatorname{dv}_{\hat{g}}=\left(\Delta_{g} \varphi-2\right) \mathrm{dv}_{g}$, we have as $\varepsilon \rightarrow 0$

$$
C_{\varepsilon}(\mathbf{x}) \rightarrow 2 \pi \sum_{i \neq j} \alpha_{i} \alpha_{j} G_{g}\left(x_{i}, x_{j}\right)+\sum_{i}\left(\frac{\alpha_{i}^{2}}{4}-\frac{Q \alpha_{i}}{2}\right) \varphi\left(x_{i}\right)+\frac{Q}{2} \sum_{i} \alpha_{i} c_{g}(\varphi)+\frac{1}{2} \sum_{i} \alpha_{i}^{2} W_{g}\left(x_{i}\right) .
$$

By applying Girsanov transform just like in the proof of Proposition 4.2, we can get rid of the $\left\langle X_{g}, K_{\hat{g}}\right\rangle_{g}$ term and this shifts the field $\hat{X}_{g, \varepsilon}$ in $Z_{\varepsilon}$ by $F(x)=-\frac{Q}{2}\left(\varphi(x)-c_{g}(\varphi)\right)$ :

$$
A_{\varepsilon}=e^{C_{\varepsilon}(\mathbf{x})+\frac{Q^{2}}{16 \pi}\|d \varphi\|_{L_{g}^{2}}^{2}} \int_{\mathbb{R}} e^{c\left(\sum_{i} \alpha_{i}-Q \chi(M)\right)} \mathbb{E}\left[\exp \left(-\mu e^{\gamma c} \hat{Z}_{\varepsilon}\right)\right] d c(1+o(1))
$$

where $\hat{Z}_{\varepsilon}:=\varepsilon^{\frac{\gamma^{2}}{2}} \int_{M} e^{\gamma\left(\hat{X}_{g, \varepsilon}+H_{g, \varepsilon}+F\right)} \mathrm{dv}_{\hat{g}}$; here we have denoted $c_{g}(\varphi)=\langle\varphi, 1\rangle_{g} / \operatorname{Vol}_{g}(M)$. By Lemma 3.2, $\left\|H_{\hat{g}, \varepsilon}\right\|_{L^{\infty}}<\infty$ thus by Proposition 3.3, we get that $\mathbb{E}\left[\hat{Z}_{\varepsilon}\right]<\infty$ thus we can find
$B>0$ such that $\mathbb{P}\left(\hat{Z}_{\varepsilon} \leqslant B\right)>0$. We therefore get

$$
A_{\varepsilon} \geqslant \beta_{\varepsilon, \mathbf{x}} \int_{-\infty}^{0} e^{c\left(\sum_{i} \alpha_{i}-Q \chi(M)\right)-\mu e^{\gamma c}} \mathbb{P}\left(\hat{Z}_{\varepsilon} \leqslant B\right) d c
$$

for some $\beta_{\varepsilon, \mathbf{x}}>0$, and this is infinite if $\sum_{i} \alpha_{i}-Q \chi(M) \leqslant 0$. Then we assume (4.9). We also have as in (3.12) the relation

$$
\hat{Z}_{\varepsilon}=e^{\frac{\gamma Q}{2} c_{g}(\varphi)} Z_{\varepsilon}(1+o(1)), \quad Z_{\varepsilon}:=\varepsilon^{\frac{\gamma^{2}}{2}} \int_{M} e^{\gamma\left(X_{g, \varepsilon}+H_{g, \varepsilon}\right)} d \mathcal{G}_{\varepsilon}^{\gamma} .
$$

Making the change of variables $c \rightarrow c-\frac{Q}{2} c_{g}(\varphi)$, we obtain that $A_{\varepsilon}$ is equal to

$$
e^{C(\mathbf{x})+\sum_{i}\left(\frac{\alpha_{i}^{2}}{4}-\frac{Q \alpha_{i}}{2}\right) \varphi\left(x_{i}\right)+\frac{Q^{2}}{16 \pi}\|d \varphi\|_{L_{g}^{2}}^{2}+\frac{Q^{2}}{2} \chi(M) c_{g}(\varphi)} \int_{\mathbb{R}} e^{c\left(\sum_{i} \alpha_{i}-Q \chi(M)\right)} \mathbb{E}\left[\exp \left(-\mu e^{\gamma c} Z_{\varepsilon}\right)\right] d c
$$

times $1+o(1)$ as $\varepsilon \rightarrow 0$, where $C(\mathbf{x})$ is given by (4.11). In particular this implies (4.13) if we can show that for the case $\varphi=0$ the limit of $A_{\varepsilon}$ is finite. We now assume $\varphi=0$, or equivalently we consider $\hat{g}=g$ the hyperbolic metric. We make the change of variable $c=\mu e^{\gamma c} Z_{\varepsilon}$ in the $c$-integral defining $A_{\varepsilon}$ (recall that $Z_{\varepsilon}>0$ almost surely), and we get

$$
A_{\varepsilon}=\gamma^{-1} e^{C(\mathbf{x})} \mu^{\frac{-\sum_{i} \alpha_{i}+Q \chi(M)}{\gamma}} \Gamma\left(\frac{\sum_{i} \alpha_{i}-Q \chi(M)}{\gamma}\right) \mathbb{E}\left[Z_{\varepsilon}^{-\frac{\sum_{i} \alpha_{i}-Q \chi(M)}{\gamma}}\right] .
$$

It remains to show that if $\alpha_{i}<Q$ for all $i$ and $s<0$, then

$$
\begin{equation*}
\lim _{\varepsilon \rightarrow 0} \mathbb{E}\left[Z_{\varepsilon}^{s}\right]=\mathbb{E}\left[Z_{0}^{s}\right] \in(0, \infty), \quad Z_{0}:=\int_{M} e^{\gamma H_{g}} d \mathcal{G}_{\hat{g}}^{\gamma} \tag{4.12}
\end{equation*}
$$

with $H_{g}:=\lim _{\varepsilon \rightarrow 0} H_{g, \varepsilon}=2 \pi \sum_{i} \alpha_{i} G_{g}\left(x_{i}, \cdot\right)$ in the $H^{-s}(M)$ topology, and that if $\alpha_{i} \geqslant Q$ for some $i$, then $\mathbb{E}\left[Z_{\varepsilon}^{s}\right] \rightarrow 0$. But this part is only a local argument and therefore Lemma 3.3. of [DKRV] applies directly. The argument goes essentially as follows. We need to determine whether the measure $\mathcal{G}_{g}^{\gamma}(d x)$ integrates the singularity $\frac{1}{d_{g}\left(x, x_{i}\right)^{\gamma \alpha_{i}}}$ in the neighborhood of $x_{i}$. Standard multifractal analysis shows that for any $\delta>0$ one can find a constant $C_{\delta}$ such that

$$
\sup _{r<1} r^{-\gamma Q+\delta} \mathcal{G}_{g}^{\gamma}\left(B_{r}\left(x_{i}\right)\right) \leqslant C_{\delta}
$$

where $B_{r}\left(x_{i}\right)$ stands for the geodesic ball of radius $r$ centered at $x_{i}$. This gives the condition $\alpha_{i}<Q$. Finally it remains to determines whether the quantity

$$
\int_{\mathbb{R}} e^{c\left(\sum_{i} \alpha_{i}-2 Q(1-\mathbf{g})\right)} \mathbb{E}\left[e^{-\mu e^{\gamma c} \mathcal{G}_{g, \mathbf{x}, \alpha}^{\gamma}(M)}\right] d c
$$

is finite. As we have seen that $\mathcal{G}_{g, \mathbf{x}, \alpha}^{\gamma}(M)$ is a well defined non trivial random variable under the condition (4.10), one may think of it as a macroscopic quantity and replace it by a constant quantity, say 1 , so as to be left with the integral

$$
\int_{\mathbb{R}} e^{c\left(\sum_{i} \alpha_{i}-2 Q(1-\mathrm{g})\right)-\mu e^{\gamma c}} d c
$$

which is straightforwardly seen to be converging if and only if (4.9) holds.
In fact the proof of the previous Proposition (adding a functional $F$ does not change anything) also shows the

Proposition 4.5. (Conformal anomaly and diffeomorphism invariance) Let $\mathbf{x}=\left(x_{1}, \ldots, x_{n}\right) \in$ $M^{n}$ and $\alpha=\left(\alpha_{1}, \ldots, \alpha_{n}\right) \in \mathbb{R}^{n}$. Let $g \in \operatorname{Met}_{-1}(M)$ be a hyperbolic metric and $\hat{g}=e^{\omega} g$ for some $\omega \in C^{\infty}(M)$. Then

$$
\begin{equation*}
\log \frac{\Pi_{\gamma, \mu}^{\mathrm{x}, \alpha}(\hat{g}, F)}{\Pi_{\gamma, \mu}^{\mathrm{X}, \alpha}\left(g, F\left(\cdot-\frac{Q}{2}\right)\right)}=\frac{1+6 Q^{2}}{96 \pi} \int_{M}\left(|d \omega|_{g}^{2}+2 K_{g} \omega\right) \mathrm{dv}_{g}+\sum_{i}\left(\frac{\alpha_{i}^{2}}{4}-\frac{Q \alpha_{i}}{2}\right) \omega\left(x_{i}\right) \tag{4.13}
\end{equation*}
$$

Let $\psi: M \rightarrow M$ be an orientation preserving diffeomorphism. Then

$$
\Pi_{\gamma, \mu}^{\mathbf{x}, \boldsymbol{\alpha}}\left(\psi^{*} g, F\right)=\Pi_{\gamma, \mu}^{\psi(\mathbf{x}), \boldsymbol{\alpha}}(g, F(\cdot \circ \psi)) .
$$

## 5. Liouville Quantum gravity

5.1. The full partition function. The partition function of Liouville quantum gravity is a weighted integral over the moduli space of the Liouville quantum field theory coupled to a Conformal Field Theory (sometimes called matter field in physics in this context). The weight of each modulus is given by some explicit functional $Z_{\text {Ghost }}(g)$ (this weight depends on the underlying surface $(M, g)$ ), called the ghost system in physics. This takes into account the factorization of the space of metrics by the action of the group of diffeomorphisms of the surface (as explained for example in [DhPh]).

Let us first recall the physics heuristics that leads to the partition function, by following [Po, DhPh, DiKa, Da]; the following discussion is not mathematically rigorous but is rather a "state of the art" in physics literature. The partition function for (Euclidean) quantum gravity in 2 D , coupled with matter, is

$$
Z=\int_{\mathcal{R}} e^{-S_{\mathrm{EH}}(g)}\left(\int e^{-S_{\mathrm{M}}\left(g, \phi_{m}\right)} D_{g} \phi_{m}\right) D g
$$

where $\mathcal{R}=\operatorname{Met}(M) / \operatorname{Diff}(M)$ is the space of Riemannian structures, the action $S_{\mathrm{EH}}(g)=$ $\mu_{0} \operatorname{Vol}_{g}(M)$ is the Einstein-Hilbert action (or gravity action) with $\mu_{0} \in \mathbb{R}$ the cosmological constant and $\phi_{m}$ are matter fields with $S_{\mathrm{M}}\left(g, \phi_{m}\right)$ being the action for matter which depends in $g$ in a conformally invariant way. Notice that, in comparison with (1.10), we got rid of the term $\int_{M} K_{g} d \mathrm{v}_{g}$ as it is a topological invariant in $2 d$ because of the Gauss-Bonnet theorem: this is an important feature of $2 d$-quantum gravity. The quantity

$$
Z_{\mathrm{M}}(g):=\int e^{-S_{\mathrm{M}}\left(g, \phi_{m}\right)} D_{g} \phi_{m}
$$

is supposed to be a CFT with central charge $\mathbf{c}_{\mathrm{M}}, D_{g} \phi_{m}$ a formal measure depending on $g$ and $D g$ is the formal Riemannian measure induced by the $L^{2}$ Riemannian metric on $\operatorname{Met}(M)$ given by (2.5) (the group $\operatorname{Diff}(M)$ acts by isometries on $\operatorname{Met}(M)$ thus the $L^{2}$-metric on $\operatorname{Met}(M)$ descends to $\mathcal{R})$. Each metric can be decomposed as $g=\psi^{*}\left(e^{\varphi} g_{\tau}\right)$ where $\tau$ is a parameter on moduli space $\mathcal{M}_{\mathbf{g}}, g_{\tau}$ is a family of metrics representing moduli space and $\psi \in \operatorname{Diff}(M)$, and the formal measure $D g$ can be accordingly decomposed as

$$
D g=Z_{\text {Ghost }}\left(e^{\varphi} g_{\tau}\right) D_{e^{\varphi} g_{\tau}} \varphi D \tau
$$

where $Z_{\text {Ghost }}$ is the ghost determinant which comes from the Jacobian of the quotient of $\operatorname{Met}(M)$ by the group of diffeomorphism $\operatorname{Diff}(M)$ (see for example [DhPh]), and given by

$$
Z_{\text {Ghost }}(g)=\left(\frac{\operatorname{det}\left(P_{g}^{*} P_{g}\right)}{\operatorname{det} J_{g}}\right)^{1 / 2}
$$

where $P_{g}, J_{g}$ are defined in Section 2.3. The ghost determinant satisfies the conformal anomaly formula (4.3) with central charge

$$
\begin{equation*}
\mathbf{c}_{\text {ghost }}=-26 \tag{5.1}
\end{equation*}
$$

Here $D \tau$ is a measure on the slice of metrics $g_{\tau}$ chosen to represent moduli space, whose value is $D \tau:=\left(\operatorname{det} J_{g_{\tau}}\right)^{1 / 2} d \tau$ with $d \tau$ being the Weil-Petersson volume form on the moduli space $\mathcal{M}_{\mathbf{g}}$ (somehow $Z_{\text {Ghost }}\left(e^{\varphi} g_{\tau}\right) D \tau$ is the quantity that makes invariant sense, as it does not depend on the matrix $J_{g}$ ). The formal measure $D_{e^{\varphi} g_{\tau}} \varphi$ should be induced by the $L^{2}$ Riemannian metric on metrics, which on the tangent space to the conformal orbit $\left[g_{\tau}\right]=\left\{e^{\varphi} g_{\tau} ; \varphi \in C^{\infty}(M)\right\}$ is given by

$$
\|f\|_{e^{\varphi} g_{\tau}}^{2}=\int_{M} \omega^{2} e^{\varphi} \mathrm{dv}_{g_{\tau}}, \quad f=\omega e^{\varphi} g_{\tau} \in T_{e^{\varphi} g_{\tau}}\left[g_{\tau}\right] .
$$

This measure depends non-linearly on $\varphi$ and it is difficult to "do the functionnal integral" for this measure, as written in [DiKa]. Therefore David and Distler-Kawai [Da, DiKa] made the "well-motivated" assumption that

$$
e^{-S_{\mathrm{EH}}(g)} Z_{\mathrm{M}}(g) D g=Z_{\mathrm{M}}\left(g_{\tau}\right) Z_{\text {Ghost }}\left(g_{\tau}\right) e^{-S_{L}\left(g_{\tau}, \varphi\right)} D \tau D_{g_{\tau}} \varphi
$$

where $S_{L}(g, \varphi)$ is the Liouville action defined by (2.2) for some parameter $Q, \gamma, \mu$ to be chosen. A formal justification of this fact was written down later in [MaMi] and [DhKu]. Invariance of the theory by choice of slice $g_{\tau}$ representing moduli space forces the partition function $\int e^{-S_{L}\left(g_{\tau}, \varphi\right)} D_{g_{\tau}} \varphi$ to be a CFT with central charge $\mathbf{c}_{\mathrm{L}}=1+6 Q^{2}$ such that the total conformal anomaly vanishes:

$$
\mathbf{c}_{\text {ghost }}+\mathbf{c}_{\mathbf{M}}+\mathbf{c}_{\mathbf{L}}=0
$$

and $Q, \gamma$ need to be related by $Q=\gamma / 2+2 / \gamma$.
Now we stop the physics parenthesis and come back to mathematics. For the matter field, we take the particular case the most studied in the physics literature, namely

$$
\begin{equation*}
Z_{\mathrm{M}}(g):=\left(\frac{\operatorname{det}^{\prime}\left(\Delta_{g}\right)}{\operatorname{Vol}_{g}(M)}\right)^{-\mathbf{c}_{\mathrm{M}} / 2} \tag{5.2}
\end{equation*}
$$

where $\mathbf{c}_{\mathbf{M}}$ is a constant in $(-\infty, 1]$. Note that this has the central charge $\mathbf{c}_{\mathrm{M}}$ by (2.10). Furthermore, there are at least two important particular cases: pure gravity where $\mathbf{c}_{\mathrm{M}}=0$ and the $2 d$ bosonic string in the case $\mathbf{c}_{\mathrm{M}}=1$. Because it is the critical situation os this approach, the latter case is especially interesting and raises serious additional difficulties. One could consider also other CFT partition functions provided that we get an expression explicit enough to determine how it behaves at the boundary of the moduli space. For each modulus $\tau \in \mathcal{M}_{\mathrm{g}}$, we can associate a hyperbolic metric $g_{\tau}$ and we will denote by $\left(g_{\tau}\right)_{\tau}$ the family of hyperbolic metrics representing the moduli space. By definition, the partition function of Liouville quantum gravity is given by the following formula:

$$
\begin{equation*}
Z:=\int_{\mathcal{M}_{\mathbf{g}}} Z_{\mathrm{Ghost}}\left(g_{\tau}\right) \times Z_{\mathrm{M}}\left(g_{\tau}\right) \times \Pi_{\gamma, \mu}\left(g_{\tau}\right) D \tau \tag{5.3}
\end{equation*}
$$

where $D \tau:=\left(\operatorname{det} J_{g_{\tau}}\right)^{1 / 2} d \tau$ with $d \tau$ the Weil-Petersson volume form on the moduli space $\mathcal{M}_{\mathbf{g}}$, and $\Pi_{\gamma, \mu}(g)$ is the partition function of the Liouville quantum field theory. This can be reduced to

$$
\begin{equation*}
Z=C_{\mathbf{g}} \int_{\mathcal{M}_{\mathbf{g}}} \operatorname{det}\left(P_{g_{\tau}}^{*} P_{g_{\tau}}\right)^{1 / 2} \times \operatorname{det}^{\prime}\left(\Delta_{g_{\tau}}\right)^{-\mathbf{c}_{\mathrm{M}} / 2} \times \Pi_{\gamma, \mu}\left(g_{\tau}\right) d \tau \tag{5.4}
\end{equation*}
$$

with $\mathrm{C}_{\mathbf{g}}$ a constant depending only on the genus of $M$. Furthermore, as we have defined LQFT only for $\gamma \in] 0,2\left[\right.$, this shows that the central charge $\mathbf{c}_{\mathbf{M}}$ of the matter field must satisfy $\mathbf{c}_{\mathbf{M}}<1$ and, combining with the expression $Q=\frac{\gamma}{2}+\frac{2}{\gamma}$, we obtain another KPZ relation [KPZ]

$$
\begin{equation*}
\gamma=\frac{\sqrt{25-\mathbf{c}_{\mathbf{M}}}-\sqrt{1-\mathbf{c}_{\mathrm{M}}}}{\sqrt{6}}, \tag{5.5}
\end{equation*}
$$

which fixes the value of $\gamma$ in terms of $\mathbf{c}_{\mathbf{M}}$.
Now, the main result of this section is the following
Theorem 5.1. If $\gamma \in] 0,2]$ and $\mathbf{c}_{\mathrm{M}}$ satisfies relation (5.5), the partition function $Z$ given by (5.3) is finite. Hence it gives rise to a finite measure $\nu$ on $\operatorname{Rad}(M) \times \mathcal{M}_{\mathbf{g}}$ defined as follows: if $\left(g_{\tau}\right)_{\tau}$ is a family of hyperbolic metrics parametrizing the moduli space $\mathcal{M}_{\mathbf{g}}$, then

$$
\nu(F):=\int_{\mathcal{M}_{\mathbf{g}} \times \mathbb{R}}\left(\frac{\operatorname{det}\left(P_{g_{\tau}}^{*} P_{g_{\tau}}\right)}{\left(\operatorname{det}^{\prime} \Delta_{g_{\tau}}\right)^{\mathbf{c}_{\mathbf{M}}+1}}\right)^{\frac{1}{2}} \mathbb{E}\left[F\left(e^{\gamma c} \mathcal{G}_{g_{\tau}}^{\gamma}(d z), \tau\right) e^{-Q \chi(M) c-\mu e^{\gamma c} \mathcal{G}_{g_{\tau}}^{\gamma}(M)}\right] d \tau d c
$$

for all continuous functionals $F: \operatorname{Rad}(M) \times \mathcal{M}_{\mathbf{g}} \rightarrow \mathbb{R}$. When renormalized by its total mass $Z=\nu(1)$, it becomes a probability measure which we call $\mathbb{P}_{\left(g_{\tau}\right)_{\tau}, \mu}$ (with expectation $\left.\mathbb{E}_{\left(g_{\tau}\right)_{\tau}, \mu}\right)$ and the couple $\left(e^{\gamma c} \mathcal{G}_{g_{\tau}}^{\gamma}(d z), \tau\right)$ becomes a random variable on $\operatorname{Rad}(M) \times \mathcal{M}_{\mathbf{g}}$, which we denote by $\left(\mathcal{L}_{\gamma}, R\right)$ and stands for the volume form of the space (called Liouville quantum gravity measure) and its modulus (called quantum modulus).

Furthermore, for all continuous functionals $F: \operatorname{Rad}(M) \times \mathcal{M}_{\mathrm{g}} \rightarrow \mathbb{R}$

$$
\begin{align*}
& \mathbb{E}_{\left(g_{\tau}\right)_{\tau}, \mu}\left[F\left(\mathcal{L}_{\gamma}(d z), R\right)\right]  \tag{5.6}\\
& =\frac{\Gamma\left(\frac{2 Q(\mathbf{g}-1)}{\gamma}\right)}{\gamma Z \mu^{\frac{2 Q(\mathbf{g}-1)}{\gamma}}} \int_{\mathcal{M}_{\mathbf{g}}}\left(\frac{\operatorname{det}\left(P_{g_{\tau}}^{*} P_{g_{\tau}}\right)}{\left(\operatorname{det}^{\prime} \Delta_{g_{\tau}}\right)^{\mathbf{c} \mathbf{M}+1}}\right)^{\frac{1}{2}} \mathbb{E}\left[F\left(\xi_{\gamma} \frac{\mathcal{G}_{g_{\gamma}}^{\gamma}(d z)}{\mathcal{G}_{g_{\tau}}^{\gamma}(M)}, \tau\right) \mathcal{G}_{g_{\tau}}^{\gamma}(M)^{\frac{Q}{\gamma} \chi(M)}\right] d \tau
\end{align*}
$$

where $\xi_{\gamma}$ is a random variable with Gamma law of density $\frac{\frac{2 Q(\mathbf{g}-1)}{\gamma}}{\Gamma\left(\frac{2 Q(\mathbf{g}-1)}{\gamma}\right)} e^{-\mu x} x^{\frac{2 Q(\mathbf{g}-1)}{\gamma}-1} \mathbf{1}_{x \geqslant 0}$ and the random modulus $R$ has density

$$
\operatorname{det}\left(P_{g_{\tau}}^{*} P_{g_{\tau}}\right)^{\frac{1}{2}}\left(\frac{\operatorname{det}^{\prime}\left(\Delta_{g_{\tau}}\right)}{\operatorname{Vol}_{g_{\tau}}(M)}\right)^{-\frac{\left(\mathbf{c}_{\mathbf{M}}+1\right)}{2}} \mathbb{E}\left[\mathcal{G}_{g_{\tau}}^{\gamma}(M)^{\frac{Q}{\gamma}} \chi(M)\right]
$$

with respect to the $d \tau$ measure.
Let us make some comment on the above result. The LQG measure depends on the family of hyperbolic metrics $\left(g_{\tau}\right)_{\tau}$ but this is not an issue since it enjoys the following invariance by reparametrization: if $\left(\psi_{\tau}\right)_{\tau}$ is a family of orientation preserving diffeomorphisms, we get the following equality for all $\tau$

$$
\begin{equation*}
\mathbb{E}_{\left(\psi_{\tau}^{*} g_{\tau}\right)_{\tau}}\left[F\left(\mathcal{L}_{\gamma} \circ \psi_{\tau}\right) \mid R=\tau\right]=\mathbb{E}_{\left(g_{\tau}\right)_{\tau}}\left[F\left(\mathcal{L}_{\gamma}\right) \mid R=\tau\right] \tag{5.7}
\end{equation*}
$$

5.2. Proof of Theorem 5.1 in the case $\gamma=2$. Our purpose is to determine the behaviour of the partition function (5.3) of LQFT (for $\gamma=2$ ) at the boundary of the moduli space with $Z_{\mathrm{M}}(g)$ defined by (5.2) and $\mathbf{c}_{\mathrm{M}}=1$. According to the relation (4.6), this amounts to showing that the integral

$$
\begin{equation*}
\int_{\mathcal{M}_{\mathbf{g}}}\left(\frac{\operatorname{det}\left(P_{g_{\tau}}^{*} P_{g_{\tau}}\right)}{\left(\operatorname{det}^{\prime} \Delta_{g_{\tau}}\right)^{\mathbf{c}_{\mathbf{M}}+1}}\right)^{\frac{1}{2}} \mathbb{E}\left[\mathcal{G}_{g_{\tau}}^{\gamma}(M)^{\frac{Q}{\gamma} \chi(M)}\right] d \tau \tag{5.8}
\end{equation*}
$$

is finite. The singularities in this integral appear at the boundary of the moduli space, namely when the surface $(M, g)$ gets close to a surface with nodes ( $M_{0}, g_{0}$ ) by pinching $n_{p}$ geodesics with respective lengths $\left(\ell_{j}\right)_{j}$ on ( $M, g$ ) (see section 2.5). According to the explicit bounds (2.12) and (2.15) for the product $\left(\frac{\operatorname{det}\left(P_{g_{\tau}}^{*} P_{g_{\tau}}\right)}{\left(\operatorname{det}^{\prime} \Delta_{g_{\tau} \tau}\right)^{\mathrm{M}^{+1}}}\right)^{\frac{1}{2}}$ and the expression for the Weyl-Petersson measure (2.9), we can give an upper bound

$$
\begin{equation*}
C\left(\prod_{j^{\prime}=1}^{n_{p}} \ell_{j^{\prime}}^{-2} \prod_{\lambda_{i}<1 / 4} \lambda_{i}^{-1}\right)\left(\prod_{j=1}^{3 \mathrm{~g}-3} \ell_{j} d \ell_{j} d \theta_{j}\right) \tag{5.9}
\end{equation*}
$$

for the quantity $\left(\frac{\operatorname{det}\left(P_{g_{\tau}}^{*} P_{g_{\tau}}\right)}{\left(\operatorname{det} \Delta_{g_{\tau}}\right)^{\mathrm{M}^{+1}}}\right)^{\frac{1}{2}} d \tau$ in the coordinate system $\left(\ell_{j}, \theta_{j}\right)_{j=1, \ldots, 3 \mathbf{g}-3}$ associated to a pant decomposition. Hence it suffices to check the integrability of the expectation $\mathbb{E}\left[\mathcal{G}_{g}^{\gamma}(M)^{\frac{Q X(M)}{\gamma}}\right]$ with respect to the measure (5.9) near $\left(M_{0}, g_{0}\right)$. It turns out that the mass of the measure $\mathcal{G}_{g}^{\gamma}(M)$ will become very large when $g$ gets close to $g_{0}$ in the pinched region of the surface $(M, g)$, making the expectation $\mathbb{E}\left[\mathcal{G}_{g}^{\gamma}(M)^{\frac{Q \chi(M)}{\gamma}}\right]$ very small. We have to quantify the rate of decay of this expectation to show integrability. Proposition 2.7 describes how the Green function behaves near the pinched geodesics. The purpose of what follows is to explain how these estimates transfer to the above expectation.

We assume ( $M_{0}, g_{0}$ ) possesses $m+1$ connected components, and we consider a neighborhood of $\left(M_{0}, g_{0}\right)$; there is a subpartition $\left\{\gamma_{1}, \ldots, \gamma_{n_{p}}\right\}$ of $M$ corresponding to the pinched geodesics. Let us denote by $\left(S_{j}\right)_{j=1, \ldots, m+1}$ the connected components of $M \backslash\left(\cup_{j^{\prime}=1}^{n_{p}} \gamma_{j^{\prime}}\right)$. Each $\gamma_{j^{\prime}}$ has a collar neighborhood denoted $\mathcal{C}_{j^{\prime}}$ (see (2.29), for simplicity of notation we remove the $g$ dependance in $\left.\mathcal{C}_{j^{\prime}}(g)\right)$. We denote by $S_{j}^{\prime}$ the set obtained by removing from $S_{j}$ all the collars

$$
S_{j}^{\prime}:=\bigcap_{j^{\prime}=1}^{n_{p}}\left(S_{j} \backslash \mathcal{C}_{j^{\prime}}\right)
$$

in such a way that $M=\cup_{j=1}^{m+1} S_{j}^{\prime} \cup_{j^{\prime}=1}^{n_{p}} \mathcal{C}_{j^{\prime}}$. Furthermore, for each $j=1, \ldots, m+1$, we define $I_{j}=\left\{j^{\prime} \in\left\{1, \ldots, n_{p}\right\} ; S_{j} \cap \mathcal{C}_{j^{\prime}} \neq \emptyset\right\}$ the set of indices $j^{\prime}$ such that the collar $\mathcal{C}_{j^{\prime}}$ encounters $S_{j}$.

We define the following quantities for $x>0$ and $1 \leqslant j^{\prime} \leqslant n_{p}$

$$
\begin{array}{rlr}
\mathcal{C}_{j^{\prime}}^{+} & :=\mathcal{C}_{j^{\prime}} \cap\{\rho \geqslant 0\} & \mathcal{C}_{j^{\prime}}^{-}:=\mathcal{C}_{j^{\prime}} \cap\{\rho \leqslant 0\} \\
\mathcal{C}_{j^{\prime}}(x)^{+} & :=\mathcal{C}_{j^{\prime}} \cap\{x \leqslant \rho\} & \mathcal{C}_{j^{\prime}}(x)^{-}:=\mathcal{C}_{j^{\prime}} \cap\{\rho \leqslant-x\},
\end{array}
$$

here $\rho: \cup_{j^{\prime}} \mathcal{C}_{j^{\prime}} \rightarrow[-1,1]$ is the function in the collars so that the metric is given by (2.37). Let us denote by $\left(\varphi_{i}\right)_{1 \leqslant i \leqslant m}$ the eigenfunctions associated to the small (non zero) eigenvalues
$\left(\lambda_{i}\right)_{1 \leqslant i \leqslant m}$ and write the Green function $G$ as $\sum_{1 \leqslant i \leqslant m} \frac{1}{\lambda_{i}} \Pi_{\lambda_{i}}+A_{g}$. Consider $X_{g}^{\prime}$ the Gaussian field with covariance $2 \pi A_{g}$, which is nothing but

$$
\begin{equation*}
X_{g}^{\prime}=X_{g}-\sum_{1 \leqslant i \leqslant m}\left(2 \pi / \lambda_{i}\right)^{1 / 2} \frac{\left\langle X_{g}, \varphi_{i}\right\rangle}{(2 \pi)^{1 / 2}} \varphi_{i} . \tag{5.10}
\end{equation*}
$$

The finite sequence $\left(\frac{\left\langle X_{g}, \varphi_{i}\right\rangle}{(2 \pi)^{1 / 2}}\right)_{1 \leqslant i \leqslant m}$ is a sequence of i.i.d. standard Gaussian random variables, namely with law $\mathcal{N}(0,1)$, which we denote $\left(a_{i}\right)_{1 \leqslant i \leqslant m}$. Furthermore, $\left(a_{i}\right)_{1 \leqslant i \leqslant m}$ and $X_{g}^{\prime}$ are independent. In case the surface $\left(M_{0}, g_{0}\right)$ is not disconnected, write $A_{g}=G$ and $X_{g}^{\prime}=X_{g}$. We introduce the random measure: $\forall A \subset M$ Borel set

$$
\mathcal{G}^{\prime}(A):=\lim _{\varepsilon \rightarrow 0}(-\ln \varepsilon)^{1 / 2} \varepsilon^{2} \int_{A} e^{2 X_{g, \varepsilon}^{\prime}} d \mathrm{v}_{g}
$$

where $X_{g, \varepsilon}^{\prime}$ is the regularization of $X_{g}^{\prime}$ as in Lemma 3.2. Notice that the convergence in probability of this measure is ensured by the convergence in probability of the same measure involving the field $X_{g}$ instead of $X_{g}^{\prime}$ (both field coincide up to an additive continuous field so that convergence of one measure is equivalent to convergence of the other one). The main technical estimate we need in the proof is the following

Lemma 5.2. Let $U_{0} \subset \mathcal{M}_{\mathbf{g}}$ be a neighborhood of some metric with nodes $g_{0}$. For any $j^{\prime}, \delta>0$, $q>0, \lambda \in[0,1]$ there exists some constant $C$ such that for all $g \in U_{0}, \forall \ell, \ell^{\prime} \geqslant \ell_{j^{\prime}}$ and $A, B>0$ and $\psi: \mathbb{R} \rightarrow \mathbb{R}$ of class $C^{1}$ such that $\psi \circ F_{j^{\prime}}(-1)=\psi \circ F_{j^{\prime}}(1)=0$

$$
\mathbb{E}\left[\left(\int_{\mathcal{C}_{j^{\prime}}} \phi e^{\psi \circ F_{j^{\prime}}} d \mathcal{G}^{\prime}\right)^{-q}\right] \leqslant C_{q} A^{-q \lambda} B^{-q(1-\lambda)}\left(\ell \ell^{\prime}\right)^{\frac{1}{2}-\delta} \exp \left(C \int_{1 \leqslant|r| \leqslant \frac{2 \pi}{\min \left(\ell, \ell^{\prime}\right)^{1-\delta}}}\left|\psi^{\prime}(r)\right|^{2} d r\right) .
$$

where the function $\phi$ is defined on the collar $\mathcal{C}_{j^{\prime}}$ by $\phi(\rho)=A \mathbf{1}_{\mathcal{C}_{j^{\prime}}(\ell)^{+}}+B \mathbf{1}_{\mathcal{C}_{j^{\prime}}\left(\ell^{\prime}\right)^{-}}$, and $F_{j^{\prime}}(\rho):=$ $\frac{2 \pi}{\ell_{j^{\prime}}} \arctan \left(\frac{\ell_{j^{\prime}}}{\rho}\right)$. The constant $C_{q}$ depends on $q$ and the mapping $q \in\left[0,+\infty\left[\mapsto C_{q}\right.\right.$ is locally bounded.

The proof of this lemma is postponed to the end of this subsection. As a direct consequence we claim

Corollary 5.3. For any $j^{\prime}, \delta>0 q>0$, there exists some constant $C$ such that for all $g \in U_{0}$ and $\ell, \ell^{\prime} \geqslant \ell_{j^{\prime}}$ and $A, B>0$

$$
\begin{array}{r}
1) \quad \mathbb{E}\left[\mathcal{G}^{\prime}\left(\mathcal{C}_{j^{\prime}}(\ell)^{+}\right)^{-q}\right]+\mathbb{E}\left[\mathcal{G}^{\prime}\left(\mathcal{C}_{j^{\prime}}(\ell)^{-}\right)^{-q}\right] \leqslant C \ell^{1 / 2-\delta} . \\
\text { 2) } \mathbb{E}\left[\left(A \mathcal{G}^{\prime}\left(\mathcal{C}_{j^{\prime}}(\ell)^{+}\right)+B \mathcal{G}^{\prime}\left(\mathcal{C}_{j^{\prime}}\left(\ell^{\prime}\right)^{-}\right)\right)^{-q}\right] \leqslant C(A B)^{-q / 2}\left(\ell \ell^{\prime}\right)^{1 / 2-\delta} .
\end{array}
$$

Now we complete the proof while considering two main situations: either the surface ( $M_{0}, g_{0}$ ) is disconnected or not.

- Case $\left(M_{0}, g_{0}\right)$ is not disconnected: In that case the measure (5.9) can be estimated from above by the measure

$$
\begin{equation*}
C\left(\prod_{j^{\prime}=1}^{n_{p}} \ell_{j^{\prime}}^{-1}\right) \prod_{j=1}^{3 \mathrm{~g}-3} d \ell_{j} d \theta_{j} . \tag{5.11}
\end{equation*}
$$

Obviously we have

$$
\mathbb{E}\left[\mathcal{G}_{g}^{\gamma}(M)^{\frac{Q \chi(M)}{\gamma}}\right] \leqslant \mathbb{E}\left[\left(\sum_{j^{\prime}=1}^{n_{p}} \mathcal{G}^{\prime}\left(\mathcal{C}_{j^{\prime}}\right)\right)^{\frac{Q \chi(M)}{\gamma}}\right] .
$$

Now we observe that the cross covariances of the field $X_{g}^{\prime}$ in the various regions $\mathcal{C}_{j^{\prime}}$ are bounded by some uniform constant, i.e. $\sup _{g \in U_{0}} \sup _{j_{1}^{\prime} \neq j_{2}^{\prime}} \sup _{x \in \mathcal{C}_{j_{1}^{\prime}}, y \in \mathcal{C}_{j_{2}^{\prime}}}\left|\mathbb{E}\left[X_{g}^{\prime}(x) X_{g}^{\prime}(y)\right]\right| \leqslant C$ (see Proposition (2.2)). Kahane's inequality [Ka, Lemma 1] then tells us that, considering independent copies $\left(\widehat{\mathcal{G}}_{j^{\prime}}^{\prime}\right)_{1 \leqslant j^{\prime} \leqslant n_{p}}$ of $\mathcal{G}^{\prime}$, the latter expectation is less than (for some irrelevant constant $C$ )

$$
C \mathbb{E}\left[\left(\sum_{j^{\prime}=1}^{n_{p}} \widehat{\mathcal{G}}_{j^{\prime}}^{\prime}\left(\mathcal{C}_{j^{\prime}}\right)\right)^{\frac{Q \chi(M)}{\gamma}}\right] .
$$

Then, we use the following elementary inequality valid for $b_{1}, \ldots, b_{n} \geqslant 0$ and $w_{1}, \ldots, w_{n} \geqslant 0$ with $\sum_{j=1}^{n} w_{j}=1$ and $q>0$

$$
\begin{equation*}
\left(\sum_{i=1}^{n} b_{i}\right)^{-q} \leqslant \prod_{i=1}^{n} b_{i}^{-w_{i} q} \tag{5.12}
\end{equation*}
$$

to deduce that the above expectation is less than $\prod_{j^{\prime}=1}^{n_{p}} \mathbb{E}\left[\left(\widehat{\mathcal{G}}_{j^{\prime}}^{\prime}\left(\mathcal{C}_{j^{\prime}}\right)\right)^{\frac{Q \chi(M)}{\gamma n_{p}}}\right]$. Combining with Corollary 5.3, we obtain

$$
\mathbb{E}\left[\mathcal{G}_{g}^{\gamma}(M)^{\frac{Q_{\chi(M)}^{\gamma}}{\gamma}}\right] \leqslant C \prod_{j^{\prime}=1}^{n_{p}} \ell_{j^{\prime}}^{1-\delta}
$$

for some arbitrary $\delta>0$. This is enough to show integrability with respect to (5.11).

- Case $\left(M_{0}, g_{0}\right)$ is disconnected: Recall that the eigenfunctions $\left(\varphi_{i}\right)_{i=1, \ldots, m}$ converge uniformly on the compact subsets of each $S_{j}$ respectively towards some fixed value denoted $v_{i j}$. From Lemma 2.8, we have the estimate in the region $|\rho|>C \ell_{j^{\prime}}$ of the cusp $\mathcal{C}_{j^{\prime}} \cap S_{j}$

$$
\begin{equation*}
\varphi_{i j}^{+} \leqslant \varphi_{i} \leqslant \varphi_{i j}^{-} \tag{5.13}
\end{equation*}
$$

where we have set

$$
\begin{align*}
& \varphi_{i j}^{+}:=v_{i j}\left(1-\mathcal{S}\left(v_{i j}\right) C \epsilon\right)|\rho|^{-\lambda_{i}+C \mathcal{S}\left(v_{i j}\right) \lambda_{i}^{2}}-C \epsilon  \tag{5.14}\\
& \varphi_{i j}^{-}:=v_{i j}\left(1+\mathcal{S}\left(v_{i j}\right) C \epsilon\right)|\rho|^{-\lambda_{i}-C \mathcal{S}\left(v_{i j}\right) \lambda_{i}^{2}}+C \epsilon \tag{5.15}
\end{align*}
$$

for some constant $C>0$; here we have denoted by $\mathcal{S}$ the function $x \in \mathbb{R} \mapsto \mathcal{S}(x):=\operatorname{sign}(x)$. Restricting the integral to the cusp regions and then using (5.13), we have the estimate conditionally on the $\left(a_{i}\right)_{1 \leqslant i \leqslant m}$

$$
\begin{aligned}
\mathbb{E}\left[\left.\mathcal{G}_{g}^{\gamma}(M)^{\frac{Q \chi(M)}{\gamma}} \right\rvert\,\left(a_{i}\right)_{1 \leqslant i \leqslant m}\right] & \leqslant \mathbb{E}\left[\left.\mathcal{G}_{g}^{\gamma}\left(\cup_{j^{\prime}} \mathcal{C}_{j^{\prime}}\right)^{\frac{Q \chi(M)}{\gamma}} \right\rvert\,\left(a_{i}\right)_{1 \leqslant i} \leqslant m\right] \\
& =\mathbb{E}\left[\left.\left(\sum_{j^{\prime}} \int_{\mathcal{C}_{j^{\prime}}} e^{\sum_{i} 2\left(2 \pi / \lambda_{i}\right)^{1 / 2} a_{i} \varphi_{i}} d \mathcal{G}^{\prime}\right)^{\frac{Q \chi(M)}{\gamma}} \right\rvert\,\left(a_{i}\right)_{1 \leqslant i \leqslant m}\right] \\
& \leqslant \mathbb{E}\left[\left.\left(\sum_{j, j^{\prime}} \int_{\mathcal{C}_{j^{\prime}} \cap S_{j}} e^{\sum_{i} 2\left(2 \pi / \lambda_{i}\right)^{1 / 2} a_{i} \varphi_{i_{j}}^{S\left(a_{i}\right)}} d \mathcal{G}^{\prime}\right)^{\frac{Q \chi(M)}{\gamma}} \right\rvert\,\left(a_{i}\right)_{1 \leqslant i \leqslant m}\right],
\end{aligned}
$$

where $\varphi_{i j}^{\mathcal{S}\left(a_{i}\right)}=\varphi_{i j}^{+}$if $a_{i} \geqslant 0$ and $\varphi_{i j}^{\mathcal{S}\left(a_{i}\right)}=\varphi_{i j}^{-}$if $a_{i}<0$. Once again we can use Kahane's convexity inequality to show that there exists a collection of mutually independent random measure $\left(\mathcal{G}_{\left(j^{\prime}\right)}^{\prime}\right)_{1 \leqslant j^{\prime} \leqslant n_{p}}$ (and independent of the $\left(a_{i}\right)_{1 \leqslant i \leqslant m}$ ) such that, for some irrelevant constant $C>0$

$$
\begin{aligned}
& \mathbb{E}\left[\left(\sum_{j, j^{\prime}} \int_{\mathcal{C}_{j^{\prime}} \cap S_{j}} e^{\left.\left.\sum_{i} 2\left(2 \pi / \lambda_{i}\right)^{1 / 2} a_{i} \varphi_{i j}^{\mathcal{S}\left(a_{i}\right)} d \mathcal{G}^{\prime}\right) \left.^{\frac{Q \chi(M)}{\gamma}} \right\rvert\,\left(a_{i}\right)_{1 \leqslant i \leqslant m}\right]}\right.\right. \\
& \quad \leqslant C \mathbb{E}\left[\left.\left(\sum_{j, j^{\prime}} \int_{\mathcal{C}_{j^{\prime}} \cap S_{j}} e^{\sum_{i} 2\left(2 \pi / \lambda_{i}\right)^{1 / 2} a_{i} \varphi_{i j}^{\mathcal{S}\left(a_{i}\right)}} d \mathcal{G}_{\left(j^{\prime}\right)}^{\prime}\right)^{\frac{Q \chi(M)}{\gamma}} \right\rvert\,\left(a_{i}\right)_{1 \leqslant i \leqslant m]}\right.
\end{aligned}
$$

Now we choose a collection of real-valued random variables $\left(r_{j}\right)_{1 \leqslant j \leqslant m+1}$ that are non negative with $\sum_{j} r_{j}=1$ and measurable with respect to the family $\left(a_{i}\right)_{1 \leqslant i \leqslant m}$. The precise choice of the family $\left(r_{j}\right)_{1 \leqslant j \leqslant m+1}$ will be made later when appropriate. Given those $\left(r_{j}\right)_{1 \leqslant j \leqslant m+1}$, we construct new weights $\left(w_{j^{\prime}}\right)_{1 \leqslant j^{\prime} \leqslant n_{p}}$ as follows. For each $j=1, \ldots, m+1$, we define $I_{j}^{+}=\left\{j^{\prime} \in\left\{1, \ldots, n_{p}\right\} ; S_{j} \cap \mathcal{C}_{j^{\prime}}^{+} \neq \emptyset\right\}$ the set of indices $j^{\prime}$ such that the collar $\mathcal{C}_{j^{\prime}}^{+}$encounters $S_{j}$ and we define similarly $I_{j}^{-}$. Also, for each $j^{\prime}=1, \ldots, n_{p}$, there exists a unique $j$ such that $\mathcal{C}_{j^{\prime}}^{+} \subset S_{j}$ and we denote by $j_{+}^{\prime}$ that index. Similarly for $j_{-}^{\prime}$. Finally for $j^{\prime}=1, \ldots, n_{p}$, we define

$$
\begin{equation*}
w_{j^{\prime}}^{+}=\frac{r_{j_{+}^{\prime}}}{\left|I_{j_{+}^{\prime}}^{+}\right|+\left|I_{j_{+}^{\prime}}^{-}\right|}, \quad w_{j^{\prime}}^{-}=\frac{r_{j_{-}^{\prime}}}{\left|I_{j_{-}^{\prime}}^{+}\right|+\left|I_{j_{-}^{\prime}}^{-}\right|} \quad \text { and } \quad w_{j^{\prime}}=w_{j^{\prime}}^{+}+w_{j^{\prime}}^{-} \tag{5.16}
\end{equation*}
$$

The first observation is that

$$
\begin{equation*}
\sum_{j^{\prime}=1, \ldots, n_{p}} w_{j^{\prime}}=1 \tag{5.17}
\end{equation*}
$$

Indeed

$$
\begin{aligned}
\sum_{j^{\prime}=1, \ldots, n_{p}} w_{j^{\prime}} & =\sum_{j=1}^{m+1} \sum_{j^{\prime} \in I_{j}^{+}} w_{j^{\prime}}^{+}+\sum_{j=1}^{m+1} \sum_{j^{\prime} \in I_{j}^{-}} w_{j^{\prime}}^{-} \\
& =\sum_{j=1}^{m+1} \sum_{j^{\prime} \in I_{j}^{+}} \frac{r_{j}}{\left|I_{j}^{+}\right|+\left|I_{j}^{-}\right|}+\sum_{j=1}^{m+1} \sum_{j^{\prime} \in I_{j}^{-}} \frac{r_{j}}{\left|I_{j}^{+}\right|+\left|I_{j}^{-}\right|} \\
& =\sum_{j=1}^{m+1} r_{j}=1
\end{aligned}
$$

Relation (5.17) allows us to use inequality (5.12). Together with independence of the measures $\left(\mathcal{G}_{\left(j^{\prime}\right)}^{\prime}\right)_{1 \leqslant j^{\prime} \leqslant n_{p}}$ conditionally on the $\left(a_{i}\right)_{1 \leqslant i \leqslant m}$, this yields

$$
\begin{aligned}
& \mathbb{E}\left[\left.\mathcal{G}_{g}^{\gamma}(M)^{\frac{Q \chi(M)}{\gamma}} \right\rvert\,\left(a_{i}\right)_{1 \leqslant i \leqslant m}\right] \\
& \leqslant C \prod_{j^{\prime}=1}^{n_{p}} \mathbb{E}\left[\left.\left(\sum_{j} \int_{\mathcal{C}_{j^{\prime}} \cap S_{j}} e^{\sum_{i} 2\left(2 \pi / \lambda_{i}\right)^{1 / 2} a_{i} \varphi_{i j}^{\mathcal{S}\left(a_{i}\right)}} d \mathcal{G}^{\prime}\right)^{w_{j^{\prime}} \frac{Q \chi(M)}{\gamma}} \right\rvert\,\left(a_{i}\right)_{1 \leqslant i \leqslant m] .}\right.
\end{aligned}
$$

Notice that the sum over $j$ in the latter expression contains at most two non trivial terms as a cusp $\mathcal{C}_{j^{\prime}}$ possesses at most two non trivial intersections with the $S_{j}$ 's. More precisely

$$
\sum_{j} \int_{\mathcal{C}_{j^{\prime}} \cap S_{j}} e^{\sum_{i} 2\left(2 \pi / \lambda_{i}\right)^{1 / 2} a_{i} \varphi_{i j}^{\mathcal{S}\left(a_{i}\right)}} d \mathcal{G}^{\prime}=\int_{\mathcal{C}_{j^{\prime}}^{-}} e^{\sum_{i} 2\left(2 \pi / \lambda_{i}\right)^{1 / 2} a_{i} \varphi_{i j_{-}^{\prime}}^{\mathcal{S}\left(a_{i}\right)}} d \mathcal{G}^{\prime}+\int_{\mathcal{C}_{j^{\prime}}^{+}} e^{\sum_{i} 2\left(2 \pi / \lambda_{i}\right)^{1 / 2} a_{i} \varphi_{i j_{+}^{\prime}}^{\mathcal{S}\left(a_{i}\right)}} d \mathcal{G}^{\prime}
$$

Now introduce $t_{j^{\prime}}^{+}=F_{j^{\prime}}^{-1}(1)$ and $t_{j^{\prime}}^{-}=F_{j^{\prime}}^{-1}(-1)$ and rewrite the above integrals as

$$
\int_{\mathcal{C}_{j^{\prime}}} \phi e^{\psi \circ F_{j^{\prime}}} d \mathcal{G}^{\prime}
$$

in view of applying Lemma 5.2, with $\ell=\ell^{\prime}=\ell_{j^{\prime}}$

$$
\begin{aligned}
& \phi=e^{\sum_{i} 2\left(2 \pi / \lambda_{i}\right)^{1 / 2} a_{i} \varphi_{i j_{+}^{\prime}}^{\mathcal{S}\left(a_{i}\right)}\left(t_{j^{\prime}}^{+}\right)} \mathbf{1}_{\mathcal{C}_{j^{\prime}}^{+}}+e^{\sum_{i} 2\left(2 \pi / \lambda_{i}\right)^{1 / 2} a_{i} \varphi_{i j_{-}^{\prime}}^{\mathcal{S}\left(a_{i}\right)}\left(t_{j^{\prime}}^{-}\right)} \mathbf{1}_{\mathcal{C}_{j^{\prime}}^{-}} \\
& \psi=\sum_{i} 2\left(2 \pi / \lambda_{i}\right)^{1 / 2} a_{i}\left(\left(\varphi_{i j_{-}^{\prime}}^{\mathcal{S}\left(a_{i}\right)} \circ F_{j^{\prime}}^{-1}-\varphi_{i j_{-}^{\prime}}^{\mathcal{S}\left(a_{i}\right)}\left(t_{j^{\prime}}^{-}\right)\right) \mathbf{1}_{\mathcal{C}_{j^{\prime}}^{-}}+\left(\varphi_{i j_{+}^{\prime}}^{\mathcal{S}\left(a_{i}\right)} \circ F_{j^{\prime}}^{-1}-\varphi_{i j_{+}^{\prime}}^{\mathcal{S}\left(a_{i}\right)}\left(t_{j^{\prime}}^{+}\right)\right) \mathbf{1}_{\mathcal{C}_{j^{\prime}}^{+}}\right)
\end{aligned}
$$

Notice that $\psi \circ F_{j^{\prime}}(1)=\psi \circ F_{j^{\prime}}(-1)=0$. Then Lemma 5.2 with $\lambda=\frac{w_{j^{\prime}}^{+}}{w_{j^{\prime}}}$ gives the estimate

$$
\begin{aligned}
& \left.\mathbb{E}\left[\left(\sum_{j} \int_{\mathcal{C}_{j^{\prime}} \cap S_{j}} e^{\sum_{i} 2\left(2 \pi / \lambda_{i}\right)^{1 / 2} a_{i} \varphi_{i j}^{\mathcal{S}\left(a_{i}\right)}} d \mathcal{G}^{\prime}\right)\right)^{w_{j^{\prime}} \frac{Q \chi(M)}{\gamma}} \right\rvert\,\left(a_{i}\right)_{1 \leqslant i \leqslant m]} \\
& \leqslant C\left(e^{\frac{Q \chi(M)}{\gamma} \sum_{i} 2\left(2 \pi / \lambda_{i}\right)^{1 / 2} a_{i}\left(w_{j^{\prime}}^{-} \varphi_{i j_{-}^{\prime}}^{\mathcal{S}\left(a_{i}\right)}\left(t_{j^{\prime}}^{-}\right)+w_{j^{\prime}}^{+} \varphi_{i j_{+}^{\prime}}^{\mathcal{S}\left(a_{i}\right)}\left(t_{j^{\prime}}^{+}\right)\right)}\right) \ell_{j^{\prime}}^{1-\delta} e^{C \sum_{i} \lambda_{i} a_{i}^{2}}
\end{aligned}
$$

 we have used the chain rule formula for derivatives combined with the following elementary estimates for $F_{j^{\prime}}^{-1}(r)=\frac{\ell_{j^{\prime}}}{\tan \frac{\ell_{j^{\prime}}}{2 \pi}}$ valid for all $r \in\left[1, \frac{2 \pi}{\ell_{j^{\prime}}^{1-\delta}}\right]$ and for some irrelevant constant $c>0$

$$
\frac{c^{-1}}{r} \leqslant F_{j^{\prime}}^{-1}(r) \leqslant \frac{c}{r} \quad \text { and }-\frac{c}{r^{2}} \leqslant\left(F_{j^{\prime}}^{-1}\right)^{\prime}(r) \leqslant-\frac{c^{-1}}{r^{2}}
$$

Combining with the expressions $(5.14)+(5.15)$ we obtain the estimate

$$
\left|\psi^{\prime}(r)\right|^{2} \leqslant C\left(\sum_{i} \lambda_{i}^{1 / 2} a_{i}|r|^{\lambda_{i}-1}\right)^{2}
$$

yielding in turn, after integrating (and recalling that $0<\lambda_{i}<1 / 4$ ),

$$
\int_{1 \leqslant|r| \leqslant \frac{2 \pi}{\ell_{j^{\prime}}^{1-\delta}}}\left|\psi^{\prime}(r)\right|^{2} d r \leqslant C \sum_{i, i^{\prime}} a_{i} a_{i^{\prime}}\left(\lambda_{i} \lambda_{i^{\prime}}\right)^{1 / 2} \leqslant C \sum_{i} \lambda_{i} a_{i}^{2}
$$

for some global constant $C$ (which may change along lines). Let us make a remark here. In the statement of Lemma 5.2, the exponent $q$ is deterministic whereas here it is here random as a measurable function of the $w_{j^{\prime}}$ (hence of the $a_{i}$ 's), namely $q=w_{j^{\prime}} \frac{Q \chi(M)}{\gamma}$. Yet conditionally on the $a_{i}$ 's the exponent can be seen as deterministic. We can thus apply Lemma 5.2. The resulting constant $C_{q}$ is therefore a measurable function of the $a_{i}$ 's. Yet, because $C_{q}$ is locally bounded as a function of $q$ and because $w_{j^{\prime}} \in[0,1]$ for all $j^{\prime}$, we can obtain an overall deterministic constant $C$ in the above inequality by taking $C=\sup _{q \in\left[0, \frac{Q \chi(M)}{\gamma}\right]} C_{q}$.

So far, we have established that

$$
\mathbb{E}\left[\mathcal{G}_{g}^{\gamma}(M)^{\frac{Q \chi(M)}{\gamma}}\right] \leqslant C \mathbb{E}\left[e^{C \sum_{i} \lambda_{i} a_{i}^{2}} e^{\frac{Q \chi(M)}{\gamma} \sum_{i, j^{\prime}} 2\left(2 \pi / \lambda_{i}\right)^{1 / 2} a_{i}\left(w_{j^{\prime}}^{-} \varphi_{i j_{-}^{\prime}}^{\mathcal{S}\left(a_{j}\right)}\left(t_{j^{\prime}}^{-}\right)+w_{j^{\prime}}^{+} \varphi_{i j_{+}^{\prime}}^{\mathcal{S}\left(a_{i}\right)}\left(t_{j^{\prime}}^{+}\right)\right)}\right] \prod_{j^{\prime}} \ell_{j^{\prime}}^{1-\delta}
$$

We can convert the sums over $j^{\prime}$ in the above expectation into sums of the type $\sum_{j} \sum_{j^{\prime} \in I_{j}^{ \pm}}$ and use the relation $\varphi_{i j_{ \pm}^{\prime}}^{\mathcal{S}\left(a_{i}\right)}\left(t_{j^{\prime}}^{ \pm}\right)=v_{i j}+\mathcal{O}\left(\lambda_{i}\right)+\mathcal{O}(\epsilon)$ for $j^{\prime} \in I_{j}^{ \pm}$(the $\mathcal{O}$ entering this relation are uniform w.r.t. the $a_{i}$ 's as easily seen from the expressions (5.14) and (5.15)) to get for some $C>0$

$$
\begin{align*}
& \mathbb{E}\left[\mathcal{G}_{g}^{\gamma}(M)^{\frac{Q \chi(M)}{\gamma}}\right]  \tag{5.18}\\
& \leqslant C \mathbb{E}\left[e^{C \sum_{i} \lambda_{i} a_{i}^{2}} e^{C \sum_{i} \frac{\left|a_{i}\right|}{\sqrt{\lambda_{i}}}\left(\mathcal{O}(\epsilon)+\mathcal{O}\left(\lambda_{i}\right)\right)} e^{\left.\frac{Q \chi(M)}{\gamma} \sum_{i, j} 2\left(2 \pi / \lambda_{i}\right)^{1 / 2} a_{i} r_{j} v_{i j}\right)}\right] \prod_{j^{\prime}} \ell_{j^{\prime}}^{1-\delta} .
\end{align*}
$$

To complete the proof we use the structure of the coefficients $\left(v_{i j}\right)_{i j}$. Recall that the vectors $v_{i} \in \mathbb{R}^{m+1}$ with components $\left(v_{i j}\right)_{j}$ form an orthonormal family for the inner product $(u, v)=\sum_{j} u_{j} v_{j} \operatorname{Vol}_{g_{0}}\left(S_{j}\right)$ and the orthogonal complement of $\operatorname{span}\left(v_{i}\right)_{i}$ is the vector 1 with all components equal to 1 . Now we define the precise values of the coefficients $r_{j}$ involved in (5.18). Set $V=\max _{i j}\left|v_{i j}\right|$ and define

$$
\begin{equation*}
r_{j}=R\left(1+\frac{1}{2 m V} \sum_{i} \mathcal{S}\left(a_{i}\right) v_{i j}\right) \operatorname{Vol}_{g_{0}}\left(S_{j}\right) \tag{5.19}
\end{equation*}
$$

where $R$ is a normalizing constant such that $\sum_{j} r_{j}=1$. Observe that $0<r_{j}<1$ for all $j$ and with this choice

$$
\forall i, \quad \sum_{j} r_{j} v_{i j}=\frac{R}{2 m V} \mathcal{S}\left(a_{i}\right) .
$$

Now we plug this relation into the estimate (5.18) to get for some $C>0$

$$
\mathbb{E}\left[\mathcal{G}_{g}^{\gamma}(M)^{\frac{Q \chi(M)}{\gamma}}\right] \leqslant C \mathbb{E}\left[e^{-C \sum_{i} \frac{\left|a_{i}\right|}{\sqrt{\lambda_{i}}}\left(1+\mathcal{O}(\epsilon)+\mathcal{O}\left(\lambda_{i}\right)\right)+C \sum_{i} \lambda_{i} a_{i}^{2}}\right] \prod_{j^{\prime}=1}^{n_{p}} \ell_{j^{\prime}}^{1-\delta} .
$$

We can choose the neighborhood of $\left(M_{0}, g_{0}\right)$ in such a way that the term $\left|\mathcal{O}(\epsilon)+\mathcal{O}\left(\lambda_{i}\right)\right|$ is less than $1 / 2$ uniformly with respect to $i,\left(a_{i}\right)_{i}$, in which case taking expectation of the above expression yields

$$
\mathbb{E}\left[\mathcal{G}_{g}^{\gamma}(M)^{\frac{Q \chi(M)}{\gamma}}\right] \leqslant C \mathbb{E}\left[e^{-C \sum_{i} \frac{\left|a_{i}\right|}{\sqrt{\lambda_{i}}}+C \sum_{i} \lambda_{i} a_{i}^{2}}\right]\left(\prod_{j^{\prime}=1}^{n_{p}} \ell_{j^{\prime}}^{1-\delta}\right) \leqslant C \prod_{i=1}^{m} \lambda_{i}^{1 / 2} \prod_{j^{\prime}=1}^{n_{p}} \ell_{j^{\prime}}^{1-\delta} .
$$

This establishes integrability with respect to the measure (5.9) by using the rough estimate (2.31) on the eigenvalues $\lambda_{j}(g)$ in terms of the lengths $\ell_{k}(g)$.
5.3. Proof of Theorem 5.1 in the case $\gamma<2$. The proof of the case $\gamma<2$ is much simpler than the case $\gamma=2$. The main reason is that the quantity $\left(\frac{\operatorname{det}\left(P_{g_{\tau}}^{*} P_{g_{\tau}}\right)}{\left(\operatorname{det}^{\prime} \Delta_{g_{\tau}}\right)^{\mathrm{M}^{+1}}}\right)^{\frac{1}{2}} d \tau$ appearing in the integral (5.8) presents a more gentle behaviour at the boundary of the moduli space due to the lower central charge $\mathbf{c}_{\mathbf{M}}<1$. More precisely (2.12) and (2.15) now gives the following estimate for this term

$$
\begin{equation*}
C\left(\prod_{j=1}^{3 \mathrm{~g}-3} \ell_{j}\right)\left(\prod_{j^{\prime}=1}^{n_{p}} \ell_{j^{\prime}}^{-\frac{5-\mathbf{c}_{\mathbf{M}}}{2}} e^{-\frac{\pi^{2}\left(1-\mathbf{c}_{\mathrm{M}}\right)}{6 \ell_{j^{\prime}}}} \prod_{\lambda_{i}<1 / 4} \lambda_{i}^{-\frac{\mathbf{c}_{\mathbf{M}}+1}{2}}\right) \prod_{j} d \ell_{j} d \theta_{j}, \tag{5.20}
\end{equation*}
$$

hence an additional exponential decay in comparison with the case $\mathbf{c}_{\mathbf{M}}=1$.

We stick to the notations introduced in section 2.4 and in the proof of Theorem 5.1 in the case $\gamma=2$. We introduce $\left(a_{i}\right)_{1 \leqslant i \leqslant m}$ i.i.d. standard Gaussian random variables and consider the following Gaussian field:

$$
Y(x)=\sum_{i=1}^{m} f_{v_{i}}(x) \frac{a_{i}}{\sqrt{\nu_{i}}}
$$

where the $\nu_{i}$ 's are defined in Theorem 2.3 and $f_{v_{i}}$ in (2.30), with $v_{i}=\left(v_{i j}\right)_{1 \leqslant j \leqslant m+1} \in \mathbb{R}^{m+1}$ from Lemma 2.5. Now, recall that on each $S_{j}^{\prime}$ the field $Y(x)$ is the constant random variable $Y_{j}=\sum_{i=1}^{m} v_{i j} \frac{a_{i}}{\sqrt{\nu_{i}}}$.

By combining Proposition 2.2 and Lemma 2.5, the Green function $G_{g}$ is such that for all $x, x^{\prime} \in \cup_{1 \leqslant j \leqslant m+1} S_{j}^{\prime}$

$$
G_{g}\left(x, x^{\prime}\right) \leqslant \sum_{i=1}^{m} \frac{f_{v_{i}}(x) f_{v_{i}}\left(x^{\prime}\right)}{\nu_{j}(g)}+A_{g}\left(x, x^{\prime}\right)+\frac{C \delta^{1 / L}}{\nu_{1}}
$$

where $C, L>0$ are global constants. Hence, if we introduce a standard normalized Gaussian variable $Z$ and an independent Gaussian field $X_{g}^{\prime}$ (living in the space of distributions) with covariance $A_{g}$, we get that for all $x, x^{\prime} \in \bigcup_{1 \leqslant j \leqslant m+1} S_{j}^{\prime}$

$$
G_{g}\left(x, x^{\prime}\right) \leqslant \mathbb{E}\left[Y(x) Y\left(x^{\prime}\right)\right]+\mathbb{E}\left[X_{g}^{\prime}(x) X_{g}^{\prime}\left(x^{\prime}\right)\right]+C \delta^{1 / L} \nu_{1}^{-1} \mathbb{E} Z^{2}
$$

in such a way that Kahane's convexity inequality [Ka] ensures that for all $q>0$

$$
\begin{aligned}
& \mathbb{E}\left[\mathcal{G}_{g}^{\gamma}\left(\cup_{1} \leqslant j \leqslant m+1 S_{j}^{\prime}\right)^{-q}\right] \\
& =\mathbb{E}\left[\left(\sum_{j=1}^{m+1} \int_{S_{j}^{\prime}} e^{\gamma X_{g}(x)-\frac{\gamma^{2}}{2} \mathbb{E}\left[X_{g}(x)^{2}\right]} e^{\frac{\gamma^{2}}{2} W_{g}(x)} \mathrm{dv}_{g}(x)\right)^{-q}\right] \\
& \leqslant \mathbb{E}\left[\left(\sum_{j=1}^{m+1} \int_{S_{j}^{\prime}} e^{\gamma Y(x)-\frac{\gamma^{2}}{2} \mathbb{E}\left[Y(x)^{2}\right]+X_{g}^{\prime}(x)-\frac{\gamma^{2}}{2} \mathbb{E}\left[X_{g}^{\prime}(x)^{2}\right]+\gamma\left(C \delta^{1 / L} / \nu_{1}\right)^{1 / 2} Z-\frac{\gamma^{2}}{2} C \delta^{1 / L} / \nu_{1}} e^{\frac{\gamma^{2}}{2} W_{g}(x)} \mathrm{dv}_{g}(x)\right)^{-q}\right] \\
& \leqslant e^{\left(q+q^{2} / 2\right) \frac{\gamma^{2}}{2} \frac{C \delta^{1 / L}}{\nu_{1}}} \mathbb{E}\left[\left(\sum_{j=1}^{m+1} \int_{S_{j}^{\prime}} e^{\gamma Y(x)-\frac{\gamma^{2}}{2} \mathbb{E}\left[Y(x)^{2}\right]+\gamma X_{g}^{\prime}(x)-\frac{\gamma^{2}}{2} \mathbb{E}\left[X_{g}^{\prime}(x)^{2}\right]} e^{\frac{\gamma^{2}}{2} W_{g}(x)} \operatorname{dv}_{g}(x)\right)^{-q}\right] .
\end{aligned}
$$

Now, by Proposition 2.2 and Lemma 2.5, there is some global constant $C>0$ such that for all $j$ and $x \in S_{j}^{\prime}$ we have

$$
\begin{equation*}
W_{g}(x) \geqslant E\left[Y(x)^{2}\right]-C \tag{5.21}
\end{equation*}
$$

so that for some other global constant $C>0$

$$
\mathbb{E}\left[\mathcal{G}_{g}^{\gamma}\left(\cup_{1} \leqslant j \leqslant m+1 S_{j}^{\prime}\right)^{-q}\right] \leqslant C e^{\frac{C \delta^{1 / L}}{\nu_{1}}} \mathbb{E}\left[\left(\sum_{j=1}^{m+1} \int_{S_{j}^{\prime}} e^{\gamma Y(x)+\gamma X_{g}^{\prime}(x)-\frac{\gamma^{2}}{2} E\left[X_{g}^{\prime}(x)^{2}\right]} d \mathrm{v}_{g}(x)\right)^{-q}\right] .
$$

Notice that

$$
\sum_{j} \operatorname{Vol}_{g}\left(S_{j}\right) Y_{j}=\sum_{j} \operatorname{Vol}_{g}\left(S_{j}\right)\left(\sum_{i} v_{i j} \frac{a_{i}}{\sqrt{\nu_{i}}}\right)=\sum_{i} \frac{a_{i}}{\sqrt{\nu_{i}}}\left(\sum_{j} \operatorname{Vol}_{g}\left(S_{j}\right) v_{i j}\right)=0
$$

since the vectors $\left(v_{i}\right)_{i}$ are orthogonal to 1 in the $\|\cdot\|_{g}$ norm (2.27). Hence there exists almost surely some $j$ such that $Y_{j} \geqslant 0$. Therefore, gathering the above considerations, we have

$$
\begin{aligned}
\mathbb{E}\left[\mathcal{G}_{g}^{\gamma}\left(\cup_{j=1}^{m+1} S_{j}^{\prime}\right)^{-q}\right] & \leqslant C e^{\frac{C \delta^{1} / L}{\nu_{1}}} \mathbb{E}\left[\left(\sum_{j=1}^{m+1} \int_{S_{j}^{\prime}} e^{\left.\left.\gamma Y(x)+\gamma X_{g}^{\prime}(x)-\frac{\gamma^{2}}{2} \mathbb{E}\left[X_{g}^{\prime}(x)^{2}\right] \operatorname{dv}_{g}(x)\right)^{-q}\right]}\right.\right. \\
& \leqslant C e^{\frac{C \delta^{1 / L}}{\nu_{1}}} \sum_{j} \mathbb{E}\left[1_{Y_{j} \geqslant 0}\left(\int_{S_{j}^{\prime}} e^{\gamma Y_{j}+\gamma X_{g}^{\prime}(x)-\frac{\gamma^{2}}{2} \mathbb{E}\left[X_{g}^{\prime}(x)^{2}\right]} \mathrm{dv}_{g}(x)\right)^{-q}\right] \\
& \leqslant C e^{\frac{C \delta^{1 / L}}{\nu_{1}}} \sum_{j} \mathbb{E}\left[\left(\int_{S_{j}^{\prime}} e^{\gamma X_{g}^{\prime}(x)-\frac{\gamma^{2}}{2} \mathbb{E}\left[X_{g}^{\prime}(x)^{2}\right]} \operatorname{dv}_{g}(x)\right)^{-q}\right] \\
& \leqslant C e^{\frac{C \delta^{1} / L}{\nu_{1}}}
\end{aligned}
$$

where in the last inequality we have used the fact that the covariance of $X_{g}^{\prime}$ can be controlled independently of the size of the small eigenvalues (this is a consequence of Proposition 2.2 and Lemma 2.5). Combining with (5.20), this estimate shows integrability with respect to the measure (5.20) by recalling that $\nu_{1} \geqslant C \ell_{1}$ for some some constant $C>0$ and by choosing $\delta$ such that $C \delta^{1 / L}<\frac{\left(1-\mathrm{c}_{\mathrm{M}}\right) \pi^{2}}{6}$.

Finally, it remains to identify the relation (5.6). Starting from the definition of $\nu$, we have

$$
\begin{aligned}
& \mathbb{E}_{\left(g_{\tau}\right)_{\tau}, \mu}\left[F\left(\mathcal{L}_{\gamma}(d z), R\right)\right]=\nu(F) / \nu(1) \\
& \quad=\frac{1}{Z} \int_{\mathcal{M}_{\mathbf{g}}} \int_{\mathbb{R}}\left(\frac{\operatorname{det}\left(P_{g_{\tau}}^{*} P_{g_{\tau}}\right)}{\left(\operatorname{det}^{\prime} \Delta_{g_{\tau}}\right)^{c_{\mathbf{M}}+1}}\right)^{\frac{1}{2}} \mathbb{E}\left[F\left(e^{\gamma c} \mathcal{G}_{g_{\tau}}^{\gamma}(d z), \tau\right) e^{-Q \chi(M) c-\mu e^{\gamma c} \mathcal{G}_{g_{\tau}}^{\gamma}(M)}\right] d \tau d c .
\end{aligned}
$$

It suffices to make the change of variables $y=e^{\gamma c} \mathcal{G}_{g_{\tau}}^{\gamma}(M)$ to get

$$
\begin{aligned}
& \mathbb{E}_{\left(g_{\tau}\right)_{\tau}, \mu}\left[F\left(\mathcal{L}_{\gamma}(d z), R\right)\right] \\
= & \frac{1}{\gamma Z} \int_{\mathcal{M}_{\mathbf{g}}}\left(\frac{\operatorname{det}\left(P_{g_{\tau}}^{*} P_{g_{\tau}}\right)}{\left(\operatorname{det}^{\prime} \Delta_{g_{\tau}}\right)^{\mathbf{c}_{\mathbf{M}}+1}}\right)^{\frac{1}{2}} \mathbb{E}\left[F\left(y \frac{\mathcal{G}_{\mathcal{G}_{\tau}}^{\gamma}(d z)}{\mathcal{G}_{g_{\tau}}^{\gamma}(M)}, \tau\right) \mathcal{G}_{g_{\tau}}^{\gamma}(M)^{-\frac{2 Q(\mathbf{g}-1)}{\gamma}}\right] y^{\frac{2 Q(\mathbf{g}-1)}{\gamma}-1} e^{-\mu y} d y,
\end{aligned}
$$

from which our claim follows.
5.4. Proof of Lemma 5.2. We will use the results in Lemma 2.42 and Proposition 2.7. So it is convenient to introduce the notations

$$
\begin{array}{ll}
\mathcal{C}_{j^{\prime}}^{ \pm}:=\mathcal{C}_{j^{\prime}}(\ell)^{+} \cup \mathcal{C}_{j^{\prime}}\left(\ell^{\prime}\right)^{-}, & F_{j^{\prime}}(\rho):=\frac{2 \pi}{\ell_{j^{\prime}}} \arctan \left(\frac{\ell_{j^{\prime}}}{\rho}\right), \\
h(\rho):=\sum_{i}|\rho|^{-2 \lambda_{i}}(1+\ln (1 /|\rho|)) . &
\end{array}
$$

We further denote by $F_{j^{\prime}}^{-1}$ the inverse of the function $F_{j^{\prime}}$. Finally, in what follows, $C$ stands for a generic irrelevant positive constant, the value of which may change along lines.

Proposition 2.7 ensures that the Green function $A_{g}$ (which is the covariance of the field $X_{g}^{\prime}$ defined in (5.10)) satisfies for $x=\rho e^{i \theta}$ and $x^{\prime}=\rho^{\prime} e^{i \theta}$ in the collar $\mathcal{C}_{j^{\prime}}$

$$
\begin{equation*}
A_{g}\left(x, x^{\prime}\right) \leqslant G_{g_{j^{\prime}}}\left(x, x^{\prime}\right)+C^{2} h(\rho) h\left(\rho^{\prime}\right) \tag{5.22}
\end{equation*}
$$

Otherwise stated, we have bounded the errors terms in Proposition 2.7 with the help of the function $h$. We stress that the function $\chi$ in Proposition 2.7 is worth 1 only for $|\rho| \leqslant 1 / 4$ so
that the above bound is rigorously valid only for $|\rho| \leqslant 1 / 4$. Yet in the following we assume that it is valid for $|\rho| \leqslant 1$ for notational convenience because it does not change the validity of the argument.

Then we need to decompose the Gaussian field with covariance function $G_{g_{j^{\prime}}}$ according to its radial/angular parts. So we consider two independent centered Gaussian fields $X^{r}$ and $X^{a}$ defined on the collar $\mathcal{C}_{j^{\prime}}=[-1,1]_{\rho} \times(\mathbb{R} \backslash \mathbb{Z})_{\theta}$ with respective covariance kernels $G^{r}$ and $G^{a}$ defined by

$$
\begin{align*}
& G^{r}\left(\rho, \theta, \rho^{\prime}, \theta^{\prime}\right):= \begin{cases}\min \left(F_{j^{\prime}}(|\rho|), F_{j^{\prime}}\left(\left|\rho^{\prime}\right|\right)\right)-\frac{\ell_{j^{\prime}}}{2 \pi^{2}} F_{j^{\prime}}(|\rho|) F_{j^{\prime}}\left(\left|\rho^{\prime}\right|\right)+C^{2} h(\rho) h\left(\rho^{\prime}\right) & \text { if } \rho \rho^{\prime}>0 \\
\frac{\ell_{j^{\prime}}}{2 \pi^{2}} F_{j^{\prime}}(|\rho|) F_{j}\left(\left|\rho^{\prime}\right|\right)+C^{2} h(\rho) h\left(\rho^{\prime}\right) & \text { if } \rho \rho^{\prime} \leqslant 0\end{cases}  \tag{5.23}\\
& G^{a}\left(\rho, \theta, \rho^{\prime}, \theta^{\prime}\right):=-\ln \mid 1-e^{-\left|F_{j^{\prime}}(\rho)-F_{j^{\prime}}\left(\rho^{\prime}\right)\right|+2 i \pi\left(\theta-\theta^{\prime}\right) \mid \mathbf{1}_{\left\{\rho \rho^{\prime} \geqslant 0\right\}} .} \tag{5.24}
\end{align*}
$$

The two quantities $G^{r}$ and $G^{a}$ are positive definite, hence the existence of such fields (pointwise for $X^{r}$ and distributional for $X^{a}$ ). Combining the above two relations we get that for $x=\rho e^{i \theta}$ and $x^{\prime}=\rho^{\prime} e^{i \theta}$ in the collar $\mathcal{C}_{j^{\prime}}$

$$
\begin{equation*}
\mathbb{E}\left[X_{g}^{\prime}(x) X_{g}^{\prime}\left(x^{\prime}\right)\right] \leqslant \mathbb{E}\left[X^{r}(x) X^{r}\left(x^{\prime}\right)\right]+\mathbb{E}\left[X^{a}(x) X^{a}\left(x^{\prime}\right)\right]+\mathbb{E}\left[\left(\delta^{(0)} h(\rho)\right)\left(\delta^{(0)} h\left(\rho^{\prime}\right)\right)\right] \tag{5.25}
\end{equation*}
$$

where $\delta^{(0)}$ is a standard Gaussian random variable independent of everything. Hence we can use Kahane's inequality [Ka, Lemma 1] to get the estimate (recall that $m_{g_{j^{\prime}}}$ is the Robin constant defined in Lemma 2.42) for $q>0$

$$
\begin{align*}
& \mathbb{E}\left[\left(\int \phi e^{\psi \circ F_{j^{\prime}}} d \mathcal{G}^{\prime}\right)^{-q}\right]  \tag{5.26}\\
& \quad \leqslant C_{q} \mathbb{E}\left[\left(\int_{\mathcal{C}_{j^{\prime}}^{ \pm}} \phi e^{\psi \circ F_{j^{\prime}}} e^{2 C \delta^{(0)} h-2 h^{2}} e^{2 X^{r}-2 \mathbb{E}\left[\left(X^{r}\right)^{2}\right]} e^{4 \pi m_{g_{j^{\prime}}}} d \mathcal{G}_{g}^{a}\right)^{-q}\right]
\end{align*}
$$

where $C_{q}$ is an explicit constant such that $\ln C_{q}$ is quadratic in $q$. Here we have defined the random measure $\mathcal{G}_{g}^{a}$ as the limit in law as $\epsilon \rightarrow 0$ (eventually up to some subsequence) of the family of random measures $(-\ln \varepsilon)^{1 / 2} e^{2 X_{\varepsilon}^{a}-2 \mathbb{E}\left[\left(X_{\varepsilon}^{a}\right)^{2}\right]} d \mathrm{v}_{g}$. From (2.42)+Proposition 2.7, we get that $4 \pi m_{g_{j^{\prime}}}-2 \mathbb{E}\left[X_{r}^{2}\right] \geqslant 2 \ln |\rho|-2 C^{2} h(\rho)^{2}$ for $|\rho|>\ell_{j^{\prime}}$ for some constant $C>0$. Hence, for $|\rho|>\ell_{j^{\prime}}$, the measure $e^{2 C \delta^{(0)} h-2 h^{2}+2 X^{r}-2 \mathbb{E}\left[\left(X^{r}\right)^{2}\right]} e^{4 \pi m_{g^{\prime}}} d \mathcal{G}_{g}^{a}$ is greater than the measure $e^{2 X^{r}} e^{g(\rho)} \rho^{-2} d \mathcal{G}_{g}^{a}$ where we have set

$$
\begin{equation*}
g(\rho):=2 C \delta^{(0)} h(\rho)+4 \ln |\rho|-4 C^{2} h(\rho)^{2} . \tag{5.27}
\end{equation*}
$$

Now that we have simplified the deterministic part of the measure we analyze the random part. For this, we write a path decomposition result for the process $X^{r}$
Lemma 5.4. Let us consider two standard Gaussian r.v. $\delta^{(1)}, \delta^{(2)} \sim \mathcal{N}(0,1)$ and two standard Brownian bridges $\left(\operatorname{Br}_{\rho}^{+}\right)_{\rho \in[0,1]}\left(\operatorname{Br}_{\rho}^{-}\right)_{\rho \in[0,1]}$, all of them mutually independent. We have the following equality in law in the sense of processes for $\rho \in[-1,1]$

Proof. Recall the covariance structure of the Brownian bridge $\mathbb{E}\left[\operatorname{Br}_{\rho}^{+} \operatorname{Br}_{\rho^{\prime}}^{+}\right]=\rho \wedge \rho^{\prime}-\rho \rho^{\prime}$. Hence for arbitrary constants $a, c>0$ and $\rho<1 / c$

$$
\mathbb{E}\left[\left(a \mathrm{Br}_{c \rho}^{+}\right)\left(a \mathrm{Br}_{c \rho^{\prime}}^{+}\right)\right]=a^{2} c \rho \wedge \rho^{\prime}-a^{2} c^{2} \rho \rho^{\prime} .
$$

Adjusting the constants $a, c$ to fit with the covariance function $\min \left(\rho, \rho^{\prime}\right)-\frac{\ell_{j^{\prime}}}{\pi^{2}} \rho \rho^{\prime}$ leads to $a=\pi / \sqrt{\ell_{j^{\prime}}}$ and $c=\frac{\ell_{j^{\prime}}}{\pi^{2}}$. One completes easily the proof of the claim by time changing with $F_{j^{\prime}}$ and adding the covariance structure of the term $\frac{\sqrt{\ell_{j^{\prime}} / 2}}{\pi} F_{j^{\prime}}(|\rho|) \delta^{(1)}+C h(\rho) \delta^{(2)}$.

Now, observe that if we restrict to those $\ell_{j^{\prime}} \leqslant|\rho| \leqslant 1$ then $\frac{\ell_{j}}{\pi^{2}} F_{j^{\prime}}(|\rho|) \in\left[0, \frac{1}{2}\right]$ and the law of the Brownian bridge Br on $\left[0, \frac{1}{2}\right]$ is absolutely continuous with respect to the law of Brownian motion $B$ with density $2 e^{-\left|B_{1 / 2}\right|^{2}}$. Using this relation with $\mathrm{Br}^{+}$and $\mathrm{Br}^{-}$together with the scaling relation for Brownian motion $a B_{t / a^{2}} \stackrel{\text { law }}{=} B_{t}$ for fixed $a>0$, we deduce using (5.26) and Lemma 5.4 that

$$
\begin{align*}
& \mathbb{E}\left[\left(\int \phi e^{\psi \circ F_{j^{\prime}}} d \mathcal{G}^{\prime}\right)^{-q}\right]  \tag{5.28}\\
& \leqslant C_{q} \mathbb{E}\left[\left(A \int_{\mathcal{C}_{j^{\prime}}(\ell)^{+}} e^{\psi \circ F_{j^{\prime}}} e^{2 B_{F_{j^{\prime}}(\rho)}^{+}+\Theta(\rho)} \rho^{-2} d \mathcal{G}_{g}^{a}+B \int_{\mathcal{C}_{j^{\prime}}\left(\ell^{\prime}\right)^{-}} e^{\psi \circ F_{j^{\prime}}} e^{2 B_{F_{j^{\prime}}(\rho)}^{-}+\Theta(\rho)} \rho^{-2} d \mathcal{G}_{g}^{a}\right)^{-q}\right],
\end{align*}
$$

where $B^{+}, B^{-}$are two independent Brownian motions, independent of everything, and the function $\Theta$ is defined by $\Theta(\rho):=\sqrt{2 \ell_{j^{\prime}}} / \pi F_{j^{\prime}}(\rho) \delta^{(1)}+2 C h(|\rho|) \delta^{(2)}+g(\rho)$. Now we would like to get rid of the drift term $\Theta$. The point is that the behaviour of $\Theta$ is rather tricky for those $|\rho|$ that are very close to (or less than) $\ell_{j^{\prime}}$ whereas the contribution of the $\Theta(\rho)$, say for $|\rho| \geqslant \ell_{j^{\prime}}^{1-\delta}$ for any $\delta>0$, turns out to be easily controlled. So we fix an arbitrary $\delta \in] 0,1[$ and remark that the expectation in the r.h.s. of (5.28) is less than the same expectation with integration restricted to $\mathcal{C}_{j^{\prime}}\left(\ell^{\prime 1-\delta}\right)^{-}$and $\mathcal{C}_{j^{\prime}}\left(\ell^{1-\delta}\right)^{+}$. Furthermore, we introduce the (random) function $Y_{\rho}$ through the relation

$$
\begin{equation*}
\forall \rho \in] 0,1], \quad \Theta(\rho)=2 \int_{F_{j^{\prime}}(1)}^{F_{j^{\prime}}(\rho)} Y_{u} d u+\Theta(1) . \tag{5.29}
\end{equation*}
$$

We set $\kappa(\ell):=F_{j^{\prime}}\left(\ell^{1-\delta}\right)=\frac{2 \pi}{\ell_{j^{\prime}}} \arctan \left(\ell_{j^{\prime}} \ell^{\delta-1}\right)$. Then the Girsanov theorem tells us that, under the probability measure

$$
R d \mathbb{P}, \quad \text { with } R:=e^{\left.\int_{\kappa(1)}^{\kappa(\ell)}\right) Y_{r} d B_{r}^{+}+\int_{\kappa(1)}^{\kappa\left(\ell^{\prime}\right)} Y_{r} d B_{r}^{-}-\frac{1}{2} \int_{\kappa(1)}^{\kappa(1)} Y_{r}^{2} d r-\frac{1}{2} \int_{\kappa(1)}^{\kappa\left(l^{\prime}\right)} Y_{r}^{2} d r}
$$

the processes $\rho \in\left[1, \ell^{1-\delta}\right] \mapsto 2 B_{F_{j^{\prime}}(\rho)}^{+}$and $\rho \in\left[1, \ell^{\prime 1-\delta}\right] \mapsto 2 B_{F_{j^{\prime}}(\rho)}^{-}$have respectively the same laws as the processes $\rho \in\left[1, \ell^{1-\delta}\right] \mapsto 2 B_{F_{j}(\rho)}^{+}+\Theta(\rho)-\Theta(1)$ and $\rho \in\left[1, \ell^{1-\delta}\right] \mapsto 2 B_{F_{j}(\rho)}^{-}+\theta(\rho)-$ $\Theta(1)$ under $\mathbb{P}$. Therefore, using the Girsanov transform in the expectation (5.28), we get

$$
\begin{aligned}
& \mathbb{E}\left[\left(\int \phi e^{\psi \circ F_{j^{\prime}}} d \mathcal{G}^{\prime}\right)^{-q}\right] \\
& \leqslant 4 \mathbb{E}\left[R e^{-q \Theta(1)}\left(A \int_{\mathcal{C}_{j^{\prime}}\left(\ell^{1-\delta}\right)^{+}} e^{\psi \circ F_{j^{\prime}}} e^{2 B_{F_{j^{\prime}}(\rho)}^{+}} \rho^{-2} d \mathcal{G}_{g}^{a}+B \int_{\mathcal{C}_{j}\left(\ell_{j^{\prime}}^{\prime 1-\delta}\right)^{-}} e^{\psi \circ F_{j^{\prime}}} e^{2 B_{F_{j^{\prime}}(\rho)}^{-}} \rho^{-2} d \mathcal{G}_{g}^{a}\right)^{-q}\right] \\
& \leqslant 4 E_{\ell, \ell^{\prime}}(q) \mathbb{E}\left[\left(A \int_{\mathcal{C}_{j^{\prime}}\left(\ell^{1-\delta}\right)^{+}} e^{\psi \circ F_{j^{\prime}}} e^{2 B_{F_{j^{\prime}}(\rho)}^{+}} \rho^{-2} d \mathcal{G}_{g}^{a}+B \int_{\mathcal{C}_{j^{\prime}}\left(\ell^{\prime 1-\delta}\right)^{-}} e^{\psi \circ F_{j^{\prime}}} e^{2 B_{F_{j^{\prime}}(\rho)}^{-}} \rho^{-2} d \mathcal{G}_{g}^{a}\right)^{-p q}\right]^{\frac{1}{p}} .
\end{aligned}
$$

where we have used the Hölder inequality to get the last inequality with $p, m, r>1$ and $\frac{1}{p}+\frac{1}{m}+\frac{1}{r}=1$ and set

$$
\begin{aligned}
E_{\ell, \ell^{\prime}}(q) & :=\mathbb{E}\left[\left(e^{\int_{\kappa(1)}^{\kappa(\ell)} Y_{u} d B_{u}^{+}+\int_{\kappa(1)}^{\kappa\left(\ell^{\prime}\right)} Y_{u} d B_{u}^{-}-\frac{1}{2} \int_{\kappa(1)}^{\kappa(\ell)} Y_{u}^{2} d u-\frac{1}{2} \int_{\kappa(1)}^{\kappa\left(l^{\prime}\right)} Y_{u}^{2} d u}\right)^{m}\right]^{1 / m} \mathbb{E}\left[e^{-q r \Theta(1)}\right]^{1 / r} \\
& =\mathbb{E}\left[e^{\frac{m^{2}-m}{2} \int_{\kappa(1)}^{\kappa(\ell)} Y_{r}^{2} d r+\frac{m^{2}-m}{2} \int_{\kappa(1)}^{\kappa\left(l^{\prime}\right)} Y_{r}^{2} d r}\right]^{1 / q} \mathbb{E}\left[e^{-q r \Theta(1)}\right]^{1 / r} .
\end{aligned}
$$

So we have got rid of the drift term $\Theta(\rho)$ with the Girsanov trick. The cost is the constant $E_{\ell, \ell^{\prime}}(q)$ but an easy computation shows that $\sup _{\ell, \ell^{\prime}} \geqslant \ell_{j} E_{\ell, \ell^{\prime}}(q)<+\infty$ for any $q>1$. This computation is left to the reader but we give a brief convincing heuristic argument. The process $Y_{u}$ is defined by (5.29). We have already explained in the proof of Theorem 5.1 that for small $\ell_{j^{\prime}}$ and for $|\rho| \geqslant \ell_{j^{\prime}}, F_{j^{\prime}}(\rho)$ behaves like $1 / \rho$. Hence $Y(\rho)$ is with good approximation given by $-\Theta^{\prime}(1 / \rho) \rho^{-2}$. Then it is readily seen that $-\Theta^{\prime}(1 / \rho) \rho^{-2}$ is a sum of terms of the type $1 / \rho,|\rho|^{-c \lambda_{i}-1}(\ln |\rho|)^{n}($ for $n=0,1,2)$ or $\ell_{j^{\prime}}^{1 / 2}$. It is then obvious to see that the square of every possible linear combination of such terms has its $\int_{\kappa(1)}^{\kappa(\ell)}$-intergal bounded by constant independently of $\ell \geqslant \ell_{j^{\prime}}$.

We can use the same argument to explicitly determine the effect of the drift term $\psi \circ F_{j^{\prime}}$. The variance of the Girsanov transform to get rid of this term is less than

$$
C \int_{1 \leqslant|r| \leqslant 2 \pi / \min \left(\ell^{1-\delta}, \ell^{\prime 1-\delta}\right)}\left|\psi^{\prime}(r)\right|^{2} d r
$$

(here we have used the fact that $\kappa(1) \geqslant 1$ and $\left.\kappa(\ell) \leqslant 2 \pi / \ell^{1-\delta}\right)$. All in all, this entails that for arbitrary $p>1$ there exists some constant $C_{p}$ such that

$$
\begin{align*}
& \mathbb{E}\left[\left(\int \phi e^{\psi \circ F_{j^{\prime}}} d \mathcal{G}^{\prime}\right)^{-q}\right] \leqslant C_{p} \exp \left(\int_{0}^{2 \pi \max \left(\ell \ell^{\prime}\right)^{\delta-1}}\left|\psi^{\prime}(r)\right|^{2} d r\right)  \tag{5.30}\\
& \quad \times \mathbb{E}\left[\left(A \int_{\mathcal{C}_{j}\left(\ell^{1-\delta}\right)^{+}} e^{\left.\left.2 B_{F_{j}(\rho)}^{+} \rho^{-2} d \mathcal{G}_{g}^{a}+B \int_{\mathcal{C}_{j}\left(\ell^{\prime} 1-\delta\right)^{-}} e^{2 B_{F_{j}(\rho)}^{-}} \rho^{-2} d \mathcal{G}_{g}^{a}\right)^{-p q}\right]^{1 / p},}\right.\right.
\end{align*}
$$

which in turn less than

$$
\begin{align*}
C_{p}(A B)^{-\frac{q}{2}} \exp & \left(\int_{0}^{2 \pi \max \left(\ell, \ell^{\prime}\right)^{\delta-1}}\left|\psi^{\prime}(r)\right|^{2} d r\right)  \tag{5.31}\\
& \times \mathbb{E}\left[\left(\int_{\mathcal{C}_{j^{\prime}}\left(\ell^{1-\delta}\right)^{+}} e^{2 B_{F^{\prime}(\rho)}^{+}} \rho^{-2} d \mathcal{G}_{g}^{a}\right)^{-p q \lambda}\left(\int_{\mathcal{C}_{j^{\prime}}\left(\ell^{\prime 1-\delta}\right)^{-}} e^{2 B_{F_{j^{\prime}}(\rho)}^{-}} \rho^{-2} d \mathcal{G}_{g}^{a}\right)^{-p q(1-\lambda)}\right]^{\frac{1}{p}}
\end{align*}
$$

after using the elementary inequality $(a+b)^{-m} \leqslant a^{-\lambda m} b^{-(1-\lambda) m}$ for $a, b \geqslant 0, \lambda \in[0,1]$ and $m>0$.

It remains to evaluate the latter expectation. So we introduce the sets for $n, k \geqslant 0$ and $\ell>0$

$$
\begin{aligned}
A_{n}^{+}(\ell) & \left.\left.=\left\{\sup _{u \in\left[\ell^{1-\delta}, 1\right]} B_{F_{j^{\prime}}(u)}^{+}-B_{F_{j^{\prime}}(1)}^{+} \in\right] n, n+1\right]\right\} \\
A_{k}^{-}\left(\ell^{\prime}\right) & \left.\left.=\left\{\sup _{u \in\left[\ell^{\prime 1-\delta, 1]}\right.} B_{F_{j^{\prime}}(u)}^{-}-B_{F_{j^{\prime}}(1)}^{-} \in\right] k, k+1\right]\right\}
\end{aligned}
$$

as well as the stopping times

$$
\left.\left.T_{n}^{+}=\left\{\inf _{u \in] 0,1]} B_{F_{j^{\prime}}(u)}^{+}-B_{F_{j^{\prime}}(1)}^{+}=n\right\} \quad T_{n}^{-}=\left\{\inf _{u \geqslant 0} B_{F_{j^{\prime}}(u)}^{-}-B_{F_{j^{\prime}}(1)}^{-} \in\right] n, n+1\right]\right\} .
$$

Partitioning the probability space according to the events $A_{n}^{+}(\ell)$ and $A_{k}^{-}\left(\ell^{\prime}\right)$ and using subadditivity of the mapping $x \in \mathbb{R}_{+} \mapsto x^{1 / p}$ for $p>1$, we get the estimate

$$
\mathbb{E}\left[\left(\int_{\mathcal{C}_{j^{\prime}}\left(\ell^{1-\delta}\right)^{+}} e^{2 B_{F_{j^{\prime}}(\rho)}^{+}} \rho^{-2} d \mathcal{G}_{g}^{a}\right)^{-p q \lambda}\left(\int_{\mathcal{C}_{j^{\prime}}\left(\ell^{\prime} 1-\delta\right)^{-}} e^{2 B_{F_{j^{\prime}}(\rho)}^{-}} \rho^{-2} d \mathcal{G}_{g}^{a}\right)^{-p q(1-\lambda)}\right]^{1 / p} \leqslant \sum_{n, n} E_{n, k}
$$

where we have set

$$
E_{n, k}:=\mathbb{E}\left[\mathbf{1}_{A_{n}^{+}(\ell)^{\prime}} \mathbf{1}_{A_{k}^{-}\left(\ell^{\prime}\right)}\left(\int_{\mathcal{C}_{j^{\prime}}\left(\ell^{1-\delta}\right)^{+}} e^{2 B_{F_{j^{\prime}}(\rho)}^{+}} \rho^{-2} d \mathcal{G}_{g}^{a}\right)^{-p q \lambda} \times\left(\int_{\mathcal{C}_{j^{\prime}}\left(\ell^{\prime} 1-\delta\right)^{-}} e^{2 B_{F_{j^{\prime}}(\rho)}^{-}} \rho^{-2} d \mathcal{G}_{g}^{a}\right)^{-p q(1-\lambda)}\right]^{\frac{1}{p}} .
$$

The idea is now the following: the fluctuations of the Brownian motion $B^{+}$over an interval of length 1 are of order 1 . Hence over the interval $\mathcal{I}_{n}^{+}:=\left[F_{j^{\prime}}\left(T_{n}^{+}\right), F_{j^{\prime}}\left(T_{n}\right)+1\right], B^{+}$is approximately equal to $n$. Put in other words, the process $B_{F_{j^{\prime}}(u)}^{+}$is worth $n$ on the interval $\left[F_{j^{\prime}}^{-1}\left(F_{j^{\prime}}\left(T_{n}^{+}\right)+\right.\right.$ 1), $T_{n}^{+}$]. Same remark for $B^{-}$. Hence $E_{n, k}$ should be estimated by

$$
\begin{align*}
E_{n, k} \leqslant & e^{-(n+k) p q} \mathbb{P}\left(A_{n}^{+}\right)^{\frac{1}{p}} \mathbb{P}\left(A_{k}^{-}\right)^{\frac{1}{p}} \\
& \times \sup _{\left.\left.x, x^{\prime} \in\right] \ell_{j}, 1\right]} \mathbb{E}\left[\left(\int_{\left\{\rho \in I^{+}(x)\right\}} \rho^{-2} d \mathcal{G}_{g}^{a}\right)^{-p q \lambda} \times\left(\int_{\left\{\rho \in I^{-}\left(x^{\prime}\right)\right\}} \rho^{-2} d \mathcal{G}_{g}^{a}\right)^{-p q(1-\lambda)}\right]^{\frac{1}{p}} \tag{5.32}
\end{align*}
$$

where for $x \in] 0,1]$, we denote $I^{+}(x):=\left[F_{j}^{-1}\left(F_{j}(x)+1\right), x\right]$ and for $x^{\prime} \in\left[-1,0\left[, I^{-}\left(x^{\prime}\right):=\right.\right.$ $\left[-x^{\prime},-F_{j}^{-1}\left(F_{j}\left(x^{\prime}\right)+1\right)\right]$. We will conclude with the two following lemmas

Lemma 5.5. For any $q>0$, we have

$$
\sup _{\ell_{j}^{\prime} \leqslant 1} \sup _{\left.x \in] \ell_{j^{\prime}}, 1\right]} \mathbb{E}\left[\left(\int_{\left\{\rho \in I^{+}(x)\right\}} \rho^{-2} d \mathcal{G}_{g}^{a}\right)^{-q}\right]<+\infty
$$

The same property for $I^{-}\left(x^{\prime}\right)$ and $x^{\prime} \in\left[-1,-\ell_{j}[\right.$.
Lemma 5.6. There is some constant $C>0$ such that for any $\delta>0, n, k \geqslant 0$ and $\ell, \ell^{\prime} \geqslant \ell_{j^{\prime}}$

$$
\mathbb{P}\left(A_{n}^{+}(\ell)\right) \leqslant C n \ell^{\frac{1}{2}(1-\delta)} \quad \mathbb{P}\left(A_{k}^{-}\left(\ell^{\prime}\right)\right) \leqslant C k \ell^{\prime \frac{1}{2}(1-\delta)}
$$

Indeed, using Hölder inequality and Lemma 5.5 we see that the expectation involved in the r.h.s. of (5.32) is less than some constant independent of everything. Hence we get the bound $E_{n, k} \leqslant C e^{-(n+k) q p} \mathbb{P}\left(A_{n}^{+}\right) \mathbb{P}\left(A_{k}^{-}\right)$. Lemma 5.6 and summability of the series $\sum_{n, k \geqslant 0} k n e^{-(n+k) q p}$ complete the argument.

The only remaining detail to fix is to show (5.32). This is an easy task as, using the independence of $B^{+}, B^{-}, \mathcal{G}_{g}^{a}$ as well as the strong Markov property of the Brownian motion, we
have

$$
\begin{aligned}
& E_{n, k} \leqslant \mathbb{E}\left[\mathbf{1}_{A_{n}^{+}(\ell) \cap A_{k}^{-}\left(\ell^{\prime}\right)} e^{-2 \lambda p q B_{F_{j^{\prime}}(1)}^{+}-2(1-\lambda) p q B_{F_{j^{\prime}}(1)}^{-}}\right. \\
& \left.\times\left(\int_{\mathcal{I}_{n}^{+}} e^{2 B_{F_{j^{\prime}}(\rho)}^{+}-2 B_{F_{j^{\prime}}(1)}^{+}} \rho^{-2} d \mathcal{G}_{g}^{a}\right)^{-p q \lambda} \times\left(\int_{\mathcal{I}_{n}^{-}} e^{2 B_{F_{j^{\prime}}(\rho)}^{-}-2 B_{F_{j^{\prime}}(1)}^{-}} \rho^{-2} d \mathcal{G}_{g}^{a}\right)^{-p q(1-\lambda)}\right]^{\frac{1}{p}} \\
& \leqslant \mathbb{E}\left[e^{-2 \lambda p q B_{F_{j^{\prime}}(1)}^{+}-2(1-\lambda) p q B_{F_{j^{\prime}}(1)}^{-}}\right]^{\frac{1}{p}} e^{-p q(n+k)}\left(\mathbb{P}\left(A_{n}^{+}(\ell)\right) \mathbb{P}\left(A_{k}^{-}\left(\ell^{\prime}\right)\right)\right)^{\frac{1}{p}} \\
& \mathbb{E}\left[e^{-p q \lambda \min _{u \in \mathcal{I}_{n}^{s}} B_{F_{j^{\prime}}(u)}^{+}-B_{F_{j} j^{\prime}(1)}^{+}}\right]^{\frac{1}{p}} \mathbb{E}\left[e^{\left.-p q(1-\lambda) \min _{u \in \mathcal{I}_{n}^{s}} B_{F_{j^{\prime}}(u)}^{-}-B_{F_{j} j^{\prime}(1)}^{-}\right]^{\frac{1}{p}}}\right. \\
& \sup _{x, x^{\prime} \in \backslash \ell_{\left.j^{\prime}, 1\right]}} \mathbb{E}\left[\left(\int_{\left\{\rho \in I^{+}(x)\right\}} \rho^{-2} d \mathcal{G}_{g}^{a}\right)^{-p q \lambda} \times\left(\int_{\left\{\rho \in I^{-}\left(x^{\prime}\right)\right\}} \rho^{-2} d \mathcal{G}_{g}^{a}\right)^{-p q(1-\lambda)}\right]^{\frac{1}{p}} .
\end{aligned}
$$

Standard estimates about the supremum of the Brownian motion over an interval of size 1 show that the $\mathbb{E}\left[e^{\left.-p q \lambda \min _{u \in \mathcal{I}_{n}^{+}} B_{F_{j}(u)}^{+}-B_{F_{j}(1)}^{+}\right]^{1 / p}}\right.$ and $\mathbb{E}\left[e^{-p q(1-\lambda) \min _{u \in \mathcal{I}_{n}^{-}} B_{F_{j}(u)}^{-}-B_{F_{j}(1)}^{-}}\right]^{1 / p}$ are bounded by some constant independent of $\ell, \ell^{\prime}$. Hence our claim.

Proof of Lemma 5.5. Recall that Gaussian multiplicative chaos at criticality (i.e. $\gamma=2$ ) possesses moments of negative order (see [DRSV1, Prop. 5]). This entails that for any $x>0$, $\mathbb{E}\left[\left(\int_{\{x \leqslant \rho \leqslant 1\}} \rho^{-2} d \mathcal{G}_{g}^{a}\right)^{-q}\right]<+\infty$. Then we observe that we have the relation $F_{j^{\prime}}^{-1}\left(F_{j^{\prime}}(x)+\right.$ 1) $>C \ell_{j^{\prime}}$ for some irrelevant constant $C$ and for all $x>\ell_{j^{\prime}}$. Therefore, for some $C$ and all $x>\ell_{j^{\prime}}$

$$
\mathbb{E}\left[\left(\int_{I^{+}(x)} \rho^{-2} d \mathcal{G}_{g}^{a}\right)^{-q}\right] \leqslant C \mathbb{E}\left[\left(\int_{I^{+}(x)}\left(\rho^{2}+\ell_{j^{\prime}}^{2}\right)^{-1} d \mathcal{G}_{g}^{a}\right)^{-q}\right]
$$

To conclude, observe that the measure $\left(\rho^{2}+\ell_{j^{\prime}}^{2}\right)^{-1} d \mathcal{G}_{g}^{a}$ is the pushforward of the measure $e^{2 X^{a}\left(F_{j^{\prime}}^{-1}(\rho)\right)-2 \mathbb{E}\left[X^{a}\left(F_{j^{\prime}}^{-1}(\rho)\right)^{2}\right]} d \rho$ (implicitly understood as the limit of a regularized sequence) under the mapping $\rho \mapsto F_{j^{\prime}}(\rho)$. The Gaussian random distribution $\rho \mapsto X^{a}\left(F_{j^{\prime}}^{-1}(\rho)\right)$ is stationary and its law does not depend on $\ell_{j^{\prime}}$ in such a way that the process $x \in \mathbb{R}_{+} \mapsto$ $\int_{x}^{x+1} e^{2 X^{a}\left(F_{j^{\prime}}^{-1}(\rho)\right)-2 \mathbb{E}\left[X^{a}\left(F_{j^{\prime}}^{-1}(\rho)\right)^{2}\right]} d \rho$ is stationary (and its law does not depend on $\ell_{j^{\prime}}$ either) and has the same law as $\int_{I(x)}\left(\rho^{2}+\ell_{j^{\prime}}^{2}\right)^{-1} d \mathcal{G}_{g}^{a}$, hence our claim.

Proof of Lemma 5.6. Recall the standard computation related to Brownian motion

$$
\mathbb{P}\left(\sup _{r \in[0, t]} B_{r} \leqslant \beta\right) \leqslant \beta t^{-1 / 2} .
$$

Therefore

$$
\mathbb{P}\left(A_{n}^{+}(\ell)\right) \leqslant \mathbb{P}\left(\sup _{u \in\left[\ell^{1-\delta}, 1\right]} B_{F_{j^{\prime}}(u)-F_{j}(1)}^{+} \leqslant n+1\right) \leqslant(n+1)\left(F_{j^{\prime}}\left(\ell^{1-\delta}\right)-F_{j^{\prime}}(1)\right)^{-1 / 2}
$$

We conclude by noticing that $\ell_{j^{\prime}}^{-1} \arctan \left(\ell_{j^{\prime}} / x\right) \geqslant C / x$ for some constant $C$ and all $x \geqslant \ell_{j^{\prime}}$. Same argument for $A_{k}^{-}\left(\ell^{\prime}\right)$.
5.5. Relation with random planar maps. The purpose of this subsection is to write a precise mathematical conjecture relating LQG to the scaling limit of large planar maps. Following Polyakov's work [Po], it was soon acknowledged by physicists that LQG should describe the scaling limit of discretized 2 d quantum gravity given by finite triangulations of a given surface, eventually coupled with a model of statistical physics (often called matter field in the physics language), see for example the classical textbook from physics [ADJ] for a review on this problem. We will describe two situations in what follows: pure gravity (no matter) or the bosonic string embedded in $D=1$ dimension.

We consider a fixed family $\left(g_{\tau}\right)_{\tau}$ of hyperbolic metrics on a compact surface $M$ (without boundary) with genus $\mathbf{g}$ as previously and the associated Liouville measure $\mathcal{L}_{\gamma}$ under $\mathbb{E}_{\left(g_{\tau}\right)_{\tau}, \mu}[\cdot]$. Let $\mathcal{T}_{N, \mathbf{g}}$ be the set of triangulations with $N$ faces with the topology of a surface of genus $\mathbf{g}$. Since these triangulations are seen up to orientation preserving homeorphisms, there are only a finite number of such triangulations. We equip $T \in \mathcal{T}_{N, \mathbf{g}}$ with a standard metric structure $h_{T}$ where each triangle is given volume $a^{2}$. The metric structure consists in gluing flat equilateral triangles: the exact definition of the metric structure is given in Les Houches lecture notes [RhVa2] in the case of the sphere and the case we consider here does not present additional difficulties for the definition. The uniformization theorem tells us that there exists a unique $\tau_{T} \in \mathcal{M}_{\mathrm{g}}$ along with an orientation preserving diffeomorphism $\psi_{T}: T \rightarrow M$ and a conformal factor $\varphi_{T}$ (with logarithmic singularities at the images of the vertices of the triangles) such that

$$
\begin{equation*}
h_{T}=\psi_{T}^{*}\left(e^{\varphi_{T}} g_{\tau_{T}}\right) \tag{5.33}
\end{equation*}
$$

Recall that in the decomposition (5.33), the functions $\varphi_{T}$ and $\psi_{T}$ are unique except if the metric $g_{\tau_{T}}$ possesses non trivial isometries. In that case, the isometry group is finite of the form $\left(\psi^{(i)}\right)_{1 \leqslant i \leqslant n}$ and starting with a decomposition (5.33) all the other decompositions of $h_{T}$ are $\left(\left(\psi^{(i)}\right)^{-1} \circ \psi_{T}\right)^{*}\left(e^{\varphi_{T} \circ \psi^{(i)}} g_{\tau_{T}}\right)$. Therefore, in the following discussion, we will suppose that the functions $\varphi_{T}$ and $\psi_{T}$ are uniquely determined by the triangulation $T$ and if this is not the case (i.e. there exists a non trivial isometry group), we replace $e^{\varphi_{T}} g_{\tau_{T}}$ in what follows by the average $\frac{1}{n} \sum_{i=1}^{n} e^{\varphi_{T} \circ \psi^{(i)}} g_{\tau_{T}}$ : these special metrics should play no role anyway as their equivalence classes are of measure 0 with respect to the Weil-Petersson volume form.

Pure gravity. It is proved in [BeCa] that the following asymptotic holds:

$$
\begin{equation*}
\left|\mathcal{T}_{N, \mathbf{g}}\right| \underset{N \rightarrow \infty}{\sim} C_{\mathcal{T}} e^{\mu_{c} N} N^{\frac{5}{2}(\mathbf{g}-1)-1} \tag{5.34}
\end{equation*}
$$

where $C_{\mathcal{T}}>0$ and $\mu_{c}>0$ are constants. The constants $C_{\mathcal{T}}, \mu_{c}$ are non universal in the sense that one can consider quadrangulations say in the place of triangulations: in this setting, the number of quadrangulations $\mathcal{Q}_{N, \mathrm{~g}}$ of size $N$ will satisfy the asymptotic $\left|\mathcal{Q}_{N, \mathrm{~g}}\right| \underset{N \rightarrow \infty}{\sim}$ $C_{\mathcal{Q}} e^{\widetilde{\mu}_{c} N} N^{\frac{5}{2}(\mathbf{g}-1)-1}$ where $C_{\mathcal{Q}}$ is different from $C_{\mathcal{T}}$ and $\widetilde{\mu}_{c}>0$ is different from $\mu_{c}$.

We set

$$
\begin{equation*}
\bar{\mu}=\mu_{c}+a^{2} \mu, \tag{5.35}
\end{equation*}
$$

where $\mu>0$ is fixed, and we consider the following random volume form on the surface $M$, defined in terms of its functional expectation

$$
\begin{equation*}
\mathbb{E}^{a}\left[F\left(\nu_{a}\right)\right]=\frac{1}{Z_{a}} \sum_{N \geqslant 1} e^{-\bar{\mu} N} \sum_{T \in \mathcal{T}_{N, \mathbf{g}}} F\left(e^{\varphi_{T}} d \mathrm{v}_{g_{\tau_{T}}}\right) \tag{5.36}
\end{equation*}
$$

for positive bounded functions $F$ where $Z_{a}$ is a normalization constant ensuring that $\mathbb{E}^{a}[\cdot]$ is the expectation of a probability measure. We denote by $\mathbb{P}^{a}$ the probability law associated to $\mathbb{E}^{a}$.

We can now state a precise mathematical conjecture:
Conjecture 1. Under $\mathbb{P}^{a}$, the random measure $\nu_{a}$ converges in law as $a \rightarrow 0$ with $\bar{\mu}$ given by (5.35) in the space of Radon measures equipped with the topology of weak convergence towards the Liouville measure $\mathcal{L}_{\gamma}$ under $\mathbb{E}_{\left(g_{\tau}\right)_{\tau}, \mu}[\cdot]$ with parameter $\gamma=\sqrt{\frac{8}{3}}$.

The fact that $\gamma=\sqrt{\frac{8}{3}}$ can be read of the total volume of space; indeed, thanks to $(5.34)$, it is easy to show that in the above asymptotic the total volume $\nu_{a}(M)$ converges to the Gamma law with density $\frac{\mu^{\frac{5}{2}(\mathbf{g}-1)}}{\Gamma\left(\frac{5}{2}(\mathrm{~g}-1)\right)} e^{-\mu x} x^{\frac{5}{2}(\mathrm{~g}-1)-1} \mathbf{1}_{x} \geqslant 0$. This law matches the law of the total volume $\xi_{\gamma}$ of $\mathcal{L}_{\gamma}$ in Theorem 5.1 for $\frac{2 Q}{\gamma}=\frac{5}{2}$, i.e. $\gamma=\sqrt{\frac{8}{3}}$.

Finally, let us mention that conjectures similar to 1 have appeared in other topologies: the sphere [DKRV], the disk [HRV] and the torus [DRV]. However, in these other topologies, the corresponding conjectures are still completely open. Let us nevertheless mention some partial progress by Curien in $[\mathrm{Cu}]$ where appealing convergence results are proven assuming a reasonable condition that has unfortunately not been proven yet.

Bosonic string. Given a triangulation $T$, let us denote by $V_{T}$ the vertex set of the dual lattice. We consider the partition function of the bosonic string on $T$ by

$$
Z(T):=\int e^{-\frac{1}{2} \sum_{v \sim v^{\prime}}\left(x_{v}-x_{v^{\prime}}\right)^{2}} \prod_{v \in V} d x_{v}
$$

where $\sim$ denotes adjacent vertices of the dual lattice. It is expected that

$$
\begin{equation*}
\sum_{T \in \mathcal{T}_{N, \mathbf{g}}} Z(T) \underset{N \rightarrow \infty}{\sim} C_{\mathcal{T}}^{\prime} e^{\mu_{c}^{\prime} N} N^{2(\mathbf{g}-1)-1} \tag{5.37}
\end{equation*}
$$

where $C_{\mathcal{T}}^{\prime}>0$ and $\mu_{c}^{\prime}>0$ are (non universal) constants. We set

$$
\begin{equation*}
\bar{\mu}=\mu_{c}^{\prime}+a^{2} \mu \tag{5.38}
\end{equation*}
$$

where $\mu>0$ is fixed, and we consider the following random volume form on the surface $M$, defined in terms of its functional expectation

$$
\begin{equation*}
\mathbb{E}^{a}\left[F\left(\nu_{a}\right)\right]=\frac{1}{Z_{a}} \sum_{N \geqslant 1} e^{-\bar{\mu} N} \sum_{T \in \mathcal{T}_{N, \mathbf{g}}} F\left(e^{\varphi_{T}} d \mathrm{v}_{g_{\tau_{T}}}\right) Z(T) \tag{5.39}
\end{equation*}
$$

for positive bounded functions $F$ where $Z_{a}$ is a normalization constant ensuring that $\mathbb{E}^{a}[\cdot]$ is the expectation of a probability measure. We denote by $\mathbb{P}^{a}$ the probability law associated to $\mathbb{E}^{a}$.

Conjecture 2. Assume $\mathbf{g}=2$. Under $\mathbb{P}^{a}$, the random measure $\nu_{a}$ converges in law as $a \rightarrow 0$ with $\bar{\mu}$ given by (5.38) in the space of Radon measures equipped with the topology of weak convergence towards the Liouville measure $\mathcal{L}_{\gamma}$ under $\mathbb{E}_{\left(g_{\tau}\right)_{\tau}, \mu}[\cdot]$ with parameter $\gamma=2$.

The reader can find much more material on $2 d$-string theory in the review [Kleb] or the lecture notes [Polch].

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E-mail address: cguillar@dma.ens.fr
DMA, U.M.R. 8553 CNRS, École Normale Superieure, 45 rue D'Ulm, 75230 Paris cedex 05, France

E-mail address: remi.rhodes@u-pem.fr
Université Paris-Est Marne la Vallée, LAMA, 5 boulevard Descartes, Champs sur Marne, 77454 Marne La Vallée, France., Research supported in part by ANR grant Liouville (ANR-15-CE40-0013)

E-mail address: Vincent.Vargas@ens.fr

DMA, U.M.R. 8553 CNRS, École Normale Superieure, 45 rue d’Ulm, 75230 Paris cedex 05, France., Partially supported by the ANR project Liouville (ANR-15-CE40-0013)


[^0]:    ${ }^{1}$ The term "critical" refers in fact to the critical dimension $D=26$ needed to get a Weyl invariant theory without quantizing the Weyl factor $e^{\omega}$ in (1.11).

[^1]:    ${ }^{2}$ Noncritical bosonic string theory is sometimes referred to as critical $D=2$ string theory, by opposition to the critical $D=26$ string theory. The two dimensions correspond to one dimension for the embedding into $\mathbb{R}$ and one dimension for the Weyl factor: in other words the Weyl factor $\omega$ in (1.11) plays the role of a hidden dimension, see the explanations in [Polch] page 121.

