ON THE BOUNDEDNESS OF EFFECTIVE POTENTIALS ARISING FROM STRING COMPACTIFICATIONS

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ABSTRACT. We study effective potentials coming from compactifications of string theory. We show that, under mild assumptions, such potentials are bounded from below in four dimensions, giving an affirmative answer to a conjecture proposed by the second author in [14]. We also derive some sufficient conditions for the existence of critical points. All proofs and mathematical hypotheses are discussed in the context of their relevance to the physics of the problem.

CONTENTS

1. Introduction	1
2. Statement of the results	3
2.1. Discussion of the hypotheses.	6
3. Proof of theorem 2.1	8
3.1. Remarks	12
4. Existence of critical points, topological considerations and some a priori estimates	13
4.1. Non-positive Yamabe invariant	16
4.2. Families of solutions	17
5. The case $d \neq 4$	19
6. Examples	22
Appendix A. Conventions and summary of formulae	25
A.1. Inequalities	26
A.2. Theorems	26
A.3. Notation and conventions	28
References	28

1. INTRODUCTION

A decade after Einstein proposed his theory of general relativity, Kaluza and Klein showed that, by postulating an extra dimension of space, one could obtain a unified theory of gravity and electromagnetism. In modern terms, one takes space-time to be 5-dimensional, an S^1 fibered over observable 4-dimensional space-time. The 5-dimensional metric can then be decomposed into a 4-dimensional metric, a metric on S^1 , and a one-form on 4-dimensional space-time. The one-form can be identified with a U(1) connection, and the 5-dimensional Einstein action becomes a 4-dimensional Einstein-Maxwell action, with additional terms depending on the metric on S^1 .

After a long period of obscurity, the ideas of Kaluza and Klein regained popularity in the 80's after the realization that supergravity and especially superstring theory make sense in more than 4

space-time dimensions. Superstring theory requires 10-dimensional space-time [31, 18], while maximally supersymmetric supergravity and its quantum version called "M theory" make sense in 11dimensional space-time. In both cases, one makes contact with standard 4-dimensional physics by postulating that the extra dimensions form a small n = 6 or 7-dimensional compact manifold M.

A primary goal of the physics work on these compactifications is to derive an effective action in 4 dimensions. This is a functional of the 4-dimensional metric and whatever additional data parametrise the extra dimensions – its metric, and the other fields of supergravity or superstring theory – taken as functions on 4 dimensional space-time. Critical points of this effective action (in the usual sense of a variational principle) correspond to critical points of the original higherdimensional supergravity or superstring action.

The simplest case is to restrict to 4-dimensional maximally symmetric space-times (Minkowski, AdS, dS) with all other fields constant in 4 dimensions. In this case, the effective action reduces to an effective potential, a functional of the metric and other fields on M. Physically, this potential is the energy of the 4-dimensional vacuum and thus considerations of stability apply – one is particularly interested in local minima of the effective potential, and one has physical arguments that the effective potential is bounded below. This idea was turned into a precise mathematical conjecture by one of the authors in [14], which we now describe.

We consider compactification on an n = D - d-dimensional compact manifold M to d-dimensional maximally symmetric space-time (Minkowski, AdS, dS), with D = 10 and d = 4 being the case of most interest. In the D-dimensional space, we consider General Relativity coupled to matter, the latter being encoded as usual in a set (possibly empty) of field strengths $F^{(p)}$, $p = 1, \ldots$ (these are curvature terms, with the standard curvature of the Yang-Mills functional being the canonical example). After compactifying we end up with an effective potential \mathcal{V} which is completely determined by quantities living in the compact n-dimensional manifold M:

$$(1.1) \quad \mathcal{V} = \frac{1}{2} \int_{M} v^{2} \Big(-R_{g} + \frac{1}{2} \sum_{p=1}^{L} |F^{(p)}|_{g}^{2} - \frac{1}{2} T_{\text{st}} \Big) dV_{g} - \frac{3}{2} \int_{M} |\nabla_{g} v|_{g}^{2} dV_{g} + \frac{1}{2} \alpha \Big(\frac{1}{G_{N}} - \int_{M} v^{2-\frac{4}{d}} dV_{g} \Big)$$

where G_N is the *d*-dimensional Newton's constant, g is a metric on M, v is a positive function (the so-called warp factor), $F^{(p)}$ are field strengths as mentioned before and $T_{\rm st}$ is a function on M which represents the non-classical sources present in superstring theory; $\alpha \in \mathbb{R}$ is a Lagrange multiplier, in the sense that its variational equation enforces a constraint (see section 2 below and [14] for a more detailed discussion).

Exploring the physics of the problem and particular examples, one is lead to the following

Conjecture 1.1. [14]: Consider a conformal class of metrics $g = e^{2\varphi}g_0$ on an n-dimensional manifold M; then the functional (1.1) evaluated at its critical points $\frac{\delta V}{\delta v} = 0$, considered as a function on the space of all conformal factors φ with fixed warped volume and volume:

$$\int_M v^{2-\frac{4}{d}} dV_g = C_1, \ \int_M dV_g = C_2$$

and all F, is bounded below.

We will prove this conjecture under mathematically precise hypotheses as theorem 2.1.

There is a close analogy to the Positive Mass theorem [41, 42, 48] and its generalisations like the Penrose inequality [30, 46]. The Positive Mass theorem states that, asymptotically flat space-times which satisfy the constraint equations of general relativity have non-negative energy. It has been generalised to space-times which are asymptotically anti-de Sitter with a given cosmological constant $\Lambda < 0$ in [11, 47], and this is very much like a lower bound $\mathcal{V} \ge \Lambda$. While the present conjecture concerns a more restricted class of space-times, in return it does not assume any asymptotic value for Λ , only that d-dimensional space-time is anti-de Sitter for some Λ .

These results also naturally relate to problems in Conformal Geometry. In particular the sign of the Yamabe invariant of M [26, 40] plays a role in our criteria for the existence of critical points (see section 4), and it would be interesting to explore possible deeper connections between the present work and the Yamabe problem.

Throughout the paper we try to explain the role played by our hypotheses and statements in the physics of the problem. Additional physics background can be found in the review [13].

2. Statement of the results

For the rest of this section we will assume that d = 4. Let us rewrite conjecture 1.1 in a form more suitable for our goals.

A positive function v is a critical point $\frac{\delta V}{\delta v} = 0$ of the functional V if and only if it satisfies¹

(2.1)
$$\Delta_g v + \left(-\frac{1}{3}R_g + \frac{1}{6}F_g - \frac{1}{6}T_{\rm st}^g\right)v = \frac{\alpha}{6}, \ v > 0,$$

for some real number α , where $F_g = \sum_{p=1}^{L} |F^{(p)}|_g^2$. T_{st} will be allowed to depend on the metric and therefore this dependence has been written explicitly. See equation (2.13) and the discussion that follows.

Let us first recall the sense in which the parameter α is a Lagrange multiplier. Since d = 4, the first constraint in conjecture 1.1 is simply $\int_M v \, dV_g = C_1$. As explained in [14], the constant C_1 is the prefactor of the 4-dimensional Einstein action (the integrated scalar curvature), which is a physically measurable constant (essentially, the inverse of Newton's constant). Thus, we write

(2.2)
$$\int_{M} v \, dV_g = \frac{1}{G_N}$$

in terms of a fixed constant G_N .² Notice that solutions to (2.1) need not automatically satisfy (2.2). However, if there exists a solution v_0 of (2.1) satisfying (2.2), then, all other solutions will automatically satisfy this condition. Indeed, notice that every solution of (2.1) is of the form

$$(2.3) v = v_0 + w,$$

where w is a solution to the homogeneous equation associated with (2.1). Since w must satisfy $\int_M w \, dV_g = 0$ (see the proof of (1) in theorem 2.1), if follows that

(2.4)
$$\int_M v \, dV_g = \int_M v_0 \, dV_g = \frac{1}{G_N}$$

Now, (2.1), being a linear equation, has the property that, a solution v with a particular value of α can be rescaled to a solution λv for another value $\lambda \alpha$. Therefore we can choose α such that (2.2) holds. Unless stated otherwise, from now on we assume that critical points are always tuned to satisfy condition (2.2).

Evaluating the functional at v and using the constraints yields

(2.5)
$$\mathcal{V} = \frac{\alpha}{4G_N}$$

¹Notice that here a sign convention opposite to that in [14] is used; see the appendix for notation and conventions.

²One could set $G_N = 1$ and thus choose the physical unit of length. Also, we have left out a conventional factor of 16π .

Because of (2.2), the dependence of \mathcal{V} on φ and v is through the Lagrange multiplier α , and we now seek to write this dependence in a more explicit fashion.

In light of (2.5), in order to prove conjecture 1.1 we need to only discuss the case where α is negative. So, in the rest of this section we make that assumption (some features of the case $\alpha \ge 0$ are discussed in sections 4 and 5). Defining u by $v = \frac{|\alpha|}{6}u$, then it is seen to satisfy

(2.6)
$$P_g u \equiv \Delta_g u + \left(-\frac{1}{3}R_g + \frac{1}{6}F_g - \frac{1}{6}T_{\rm st}^g\right)u = -1.$$

Using (2.2) to express α in terms of u, one has

(2.7)
$$|\alpha| = \frac{6}{G_N \int_M u dV_g}$$

Writing a general metric g in the conformal class in terms of a fixed background metric g_0 , $g = e^{2\varphi}g_0$, and using (2.3), (2.4), and (2.7) into (2.5) finally gives

(2.8)
$$\mathcal{V} = -\frac{6}{4G_N^2 \int_M e^{n\varphi} u \, dV_{g_0}}.$$

By the solution of the Yamabe problem [3, 36, 43, 49], the metric g_0 can be assumed to have constant scalar curvature R_{g_0} , and henceforth we do so. This will be positive, negative, or zero according to the sign of the Yamabe invariant³ of (M, g_0) . It will be convenient to express all quantities in terms of this fixed background metric.

Under $g = e^{2\varphi}g_0$,

(2.9)
$$R_g = e^{-2\varphi} \Big(-2(n-1)\Delta_{g_0}\varphi - (n-1)(n-2)|\nabla_{g_0}\varphi|_{g_0}^2 + R_{g_0} \Big),$$

and

(2.10)
$$\Delta_g u = e^{-2\varphi} \left(\Delta_{g_0} u + (n-2) \langle \nabla_{g_0} \varphi, \nabla_{g_0} u \rangle_{g_0} \right).$$

Recall that, if $F^{(p)}$ is a *p*-form, then

(2.11)
$$|F^{(p)}|_{g}^{2} = g^{\mu_{1}\nu_{1}} \cdots g^{\mu_{p}\nu_{p}} F_{\mu_{1}\cdots\mu_{p}} F_{\nu_{1}\cdots\nu_{p}}$$
$$= e^{-2p\varphi} g_{0}^{\mu_{1}\nu_{1}} \cdots g_{0}^{\mu_{p}\nu_{p}} F_{\mu_{1}\cdots\mu_{p}} F_{\nu_{1}\cdots\nu_{p}}$$
$$= e^{-2p\varphi} |F^{(p)}|_{g_{0}}^{2}.$$

Hence, the gauge fields expressed in terms of the metric g_0 and φ become

(2.12)
$$\sum_{p=1}^{L} |F^{(p)}|_{g}^{2} = \sum_{p=1}^{L} e^{-2p\varphi} |F^{(p)}|_{g_{0}}^{2}$$

Now we need to ask how T_{st}^g transforms under $g = e^{2\varphi}g_0$. The basic example of T_{st}^g in string theory is the so-called orientifold plane, which is supported on a submanifold [8, 17]. Another common example is a Chern-Simons or topological term, as used in [7, 17]. What is important for our problem is its dependence on the conformal factor. Thus, we will take T_{st}^g to be a function or a generalised function on M, and postulate a homogeneous dependence on the conformal factor,

(2.13)
$$T_{\rm st}^g = e^{-2\beta\varphi}T_{\rm st}^{g_0},$$

which is consistent with (2.12). In the main result of this paper, theorem 2.1, we will consider different situations which will allow for different choices of β . The first two cases, $\beta = 0$ and $\beta = 1$,

³We use the term "Yamabe invariant" to denote the invariant of the conformal class, whereas the supremum over all conformal classes is called "Topological Yamabe invariant".

will be treated together, and in fact the proof in these cases works for any $0 \le \beta \le 1$. Another case of special interest is $\beta = n/2$ which is appropriate for a delta function source. Finally, a variation on this condition which includes orientifolds and many other cases is to make the simple assumption

(2.14)
$$\int_{M} T_{\rm st}^{g} \, dV_{g} = \int_{M} T_{\rm st}^{g_{0}} \, dV_{g_{0}}, \text{ for all } g = e^{2\varphi} g_{0}.$$

See sections 2.1 and 6 for a more detailed discussion on the relevance of each hypothesis for the physics of string compactifications.

Upon combining (2.9), (2.10), (2.12) and (2.13), equation (2.6) becomes

(2.15)
$$M_{g_0}u = \Delta_{g_0}u + (n-2)\langle \nabla_{g_0}\varphi, \nabla_{g_0}u \rangle_{g_0} + Uu = -e^{2\varphi},$$

with

(2.16)
$$U \equiv \frac{2}{3}(n-1)\Delta_{g_0}\varphi + \frac{1}{3}(n-1)(n-2)|\nabla_{g_0}\varphi|_{g_0}^2 - \frac{1}{3}R_{g_0} + \mathcal{F}(\varphi) - \frac{1}{6}e^{2(1-\beta)\varphi}T_{\rm st}^{g_0}$$

and

(2.17)
$$\mathcal{F}(\varphi) = \frac{1}{6} \sum_{p=1}^{L} e^{2(1-p)\varphi} |F^{(p)}|_{g_0}^2.$$

We will not need the specific form of \mathcal{F} , it being enough for our proof to notice that $\mathcal{F} \geq 0.^4$

Let us now restate conjecture 1.1 as a theorem that will be proven in the subsequent sections. After stating the results we discuss the physical meaning of the hypotheses some consequences.

Theorem 2.1. Let (M, g_0) be a compact orientable n-dimensional Riemannian manifold without boundary, let $A \in \mathbb{R}_+$. Fix a collection of smooth p forms $\{F^{(p)}\}_{p=1}^L$, and, a function $T_{st}^{g_0} \in C^{\infty}(M)$ which transforms under conformal changes as (2.13), for some $\beta \geq 0$. For any smooth function φ on M define $\mathcal{F}(\varphi)$ as in (2.17). Define $\mathcal{A} = \{\varphi \in C^{\infty}(M) \mid \int_M e^{n\varphi} dV_{g_0} = A\}$, and let $\mathcal{S} \subset \mathcal{A}$ be the set of $\varphi \in \mathcal{A}$ such that equation (2.15) has a positive solution $u = u(\varphi, g_0, T_{st}^{g_0}, \mathcal{F})$. Then:

(1) The map $\mathfrak{F}_{g_0,T^{g_0}_{st},\mathcal{F}}: \mathcal{S} \to \mathbb{R}$ given by

(2.18)
$$\mathfrak{F}_{g_0,T^{g_0}_{st},\mathcal{F}}(\varphi) = -\frac{6}{4G_N^2 \int_M e^{n\varphi} u \, dV_{g_0}}$$

is well defined.

(2) Fix $\eta > 0$ and define

$$S_{\eta} = \Big\{ \varphi \in S \Big| \int_{M} R_g \, dV_g \leq \eta \,, \text{ where } g = e^{2\varphi} g_0 \Big\}.$$

Suppose that $0 \leq \beta \leq 1$, or $\beta = \frac{n}{2}$, or (2.14) holds. Then there exists a constant $K_{\eta} \in \mathbb{R}$ such that (2.19) $\inf_{\varphi \in S_n} \mathfrak{F}_{g_0, T^{g_0}_{st}, \mathcal{F}}(\varphi) \geq K_{\eta}.$

In particular, $\mathfrak{F}_{g_0,T_{st}^{g_0},\mathcal{F}}$ is bounded from below if the scalar curvature is uniformly bounded. Moreover, if dim M = 2 then

(2.20)
$$\inf_{\varphi \in \mathcal{S}} \mathfrak{F}(\varphi) \ge K$$

⁴ This is actually not manifest in supergravity as there are Chern-Simons and other non-quadratic terms, but it is shown in Ref. [15].

for some $K \in \mathbb{R}$.

(3) Furthermore, if we define

(2.21)
$$\widetilde{\mathcal{S}}_{\eta} = \Big\{ \varphi \in \mathcal{S}_{\eta} \text{ such that } \int_{M} \left(F_{g} - T_{st}^{g} \right) dV_{g} \ge -\eta \,, \text{ where } g = e^{2\varphi} g_{0} \Big\},$$

then there exists a constant $K_{\eta} \in \mathbb{R}$ such that

(2.22)
$$\inf_{\varphi \in \widetilde{\mathcal{S}}_{\eta}} \mathfrak{F}_{g_0, T_{st}^{g_0}, \mathcal{F}}(\varphi) \ge K_{\eta}$$

for any value of β in (2.13).

A different approach to the theorem would be to impose conditions directly on the functions φ appearing in the conformal factor $e^{2\varphi}$. In this regard we prove:

Proposition 2.2. Assume the same hypotheses of Theorem 2.1. Suppose that $0 \leq \beta \leq 1$ and that $\dim(M) \geq 3$. If $B_{\eta}(0)$ is the ball of radius η in $L^{1}(M, g_{0})$ then there exists a constant $K_{\eta} \in \mathbb{R}$ such that

(2.23)
$$\inf_{\varphi \in \mathcal{S} \cap B_{\eta}(0)} \mathfrak{F}_{g_0, T_{st}^{g_0}, \mathcal{F}}(\varphi) \ge K_{\eta}.$$

2.1. Discussion of the hypotheses. Let us make some comments on the hypotheses and the content of theorem 2.1. We first discuss the physics of (2.13) and (2.14).

The "non-classical" source terms $T_{\rm st}$ differ from the usual stress-energy tensor of general relativity in two ways. First, they are generally associated to quantum effects and anomalies in string theory and M theory, as in the anomaly cancellation terms of the heterotic string called on in [8]. Second, and more importantly for us, they violate the positive energy condition. This allows finding compactifications to Minkowski space-time even in the presence of other positive contributions to the energy, again as first seen in [8]. But it means that any positive (or bounded below) energy theorem in string/M theory will require placing some condition on $T_{\rm st}$, which replaces the positive energy condition. Although the physical consistency of string/M theory implies that some such condition exists, and one can see some of its features in examples, at present no precise and sufficiently general statement of the condition has been proposed. Thus we put forward (2.13) and (2.14) as candidates, which suffice to prove our main results.

Some of the β values of interest are the "topological" case referred to above, namely, $\beta = 0$, and $\beta = \frac{n}{2}$, which is how T_{st}^g transforms when it is a delta function, since

$$1 = \int_M \delta_g \, dV_g = \int_M e^{-n\varphi} \delta_{g_0} e^{n\varphi} \, dV_{g_0} = \int_M \delta_{g_0} \, dV_{g_0}.$$

Notice that $\beta = \frac{n}{2}$ is also the case where T_{st}^g scales in the same way as in (2.11) with $p = \frac{n}{2}$, which is when the gauge fields have a conformally invariant stress tensor. Another natural case is the "naive" choice $\beta = 1$, i.e., declaring that T_{st}^g transforms in the same way as the metric.

Yet another important example is the orientifold plane mentioned above. In this case, T_{st}^g is supported on a closed submanifold $N \subset M$ in the sense that

(2.24)
$$T_{\rm st}^g \equiv 0 \text{ on } M \setminus N, \text{ and } \int_M T_{\rm st}^g \, dV_g = \int_N T_{\rm st}^g \, d\Sigma_g = 1,$$

where $d\Sigma_g$ is the induced volume element on N. The simplest way to model (2.24) is to have $T_{\rm st}^g = \delta_g(N)$, i.e., the measure assigning one to N and zero to $M \setminus N$, although choices like $T_{\rm st}^g = \chi_N$,

or some smooth approximation of it, can also be considered (χ_E being the characteristic function of a set E). If N is k-dimensional, then $d\Sigma_g = e^{k\varphi} d\Sigma_{g_0}$, and hence by analogy with the delta function, we require that in the orientifold case the string term transforms as

$$T_{\rm st}^g = e^{-k\varphi} T_{\rm st}^{g_0}.$$

This corresponds to $\beta = \frac{k}{2}$, where k is the dimension of the submanifold supporting $T_{\text{st}}^{g_0}$. Notice that $\beta = \frac{n}{2}$ and an orientifold-like string term can both be encoded in the assumption (2.14).

Finally, even if neither of these conditions were to hold, we can prove the result under the global condition (2.21), which states that, any negative contribution of $T_{\rm st}^g$ is compensated by the energy of matter and fluxes, as comes out of many analyses and discussions [17, 29, 50].

Turning to mathematical questions, the first and most obvious is whether the set S is empty or not. We address this question in section 4, where we provide conditions for the existence of solutions to equation (2.15). We will also see reasons to think that $S \neq A$ in general. Presumably this is because string/M theory effects cannot be neglected in these cases, but not much is known about this.

That the map $\mathfrak{F}_{g_0,T^{g_0}_{\mathrm{st}},\mathcal{F}}$ is well-defined is not a surprise – by Eq. (2.4), $\mathfrak{F}_{g_0,T^{g_0}_{\mathrm{st}},\mathcal{F}}$ does not depend on elements in the kernel of P_q .

In (2.19), the infimum is taken not only over all conformal factors satisfying the original constraint $\int_M e^{n\varphi} dV_{g_0} = A$, but also obeying the extra integral bound $\int_M R_g dV_g \leq \eta$.

Curvature bounds are physically appropriate whenever one studies quantum gravity and string/M theory in the language of general relativity. These theories have a preferred scale of length below which a description in terms of classical space-time breaks down, the Planck length L_{pl} for quantum gravity and M theory, and the string length L_s for string theory. Let $L = \max(L_{pl}, L_s)$, then a description in terms of a space-time metric satisfying Einstein's equations will generally only be valid when the curvature of the metric is much less than $1/L^2$. This would apply to every component of the curvature as well as its derivatives, and the *p*-form field strengths $F^{(p)}$ and their derivatives. In this sense, the physics naturally places stronger conditions, such as uniform curvature bounds⁵. As we explain below, the proof with these conditions is an elementary consequence of the more general assumptions that are adopted here.

Mathematically, since all constraints in conjecture 1.1 are integral conditions, one would prefer, even if merely for aesthetic reasons, to have any additional conditions to be integral as well, as in theorem 2.1.

The bound on the scalar curvature is also a very reasonable geometrical assumption, in the sense that curvature bounds are commonplace in Riemannian Geometry (see e.g. [4, 10, 40]). Finally, we point out that $\int_M R_g \, dV_g \leq \eta$ is a natural generalization of the situation in 2 dimensions, where the condition is automatically satisfied due the the Gauss-Bonnet formula.

We should also make a comment about our assumptions of regularity. In order to avoid technicalities that would obfuscate the main ideas, all our functions are assumed to be smooth. This obviously excludes cases such as $T_{st}^{g_0} \sim \delta_{g_0}(p)$, for some $p \in M$, or $T_{st}^{g_0} \sim \delta_{g_0}(N)$, for some submanifold N. However it will be clear that, after a suitable interpretation in the context of linear equations involving a generalized function, the same argument works (see section 6). Other regularity assumptions can also be greatly relaxed. See for example [19, 24, 44, 45] for generalizations of the techniques here employed to conditions of less regularity.

 $^{^{5}}$ We recall that curvature bounds have been successfully employed to study the long-time existence of Einstein equations [2, 28].

DISCONZI, DOUGLAS, AND PINGALI

Finally, let us address the $L^1(M, g_0)$ condition in proposition 2.2. A bound on the $L^1(M, g_0)$ norm of φ seems to be a fairly standard hypothesis if we focus exclusively on the analytical aspects of the problem. From a more geometric perspective, it can be interpreted as follows. Even smooth conformal factors can approach distributions with very bad singularities, in which case the limit metrics would be highly degenerate. On the other hand, singularities do occur in quantum field theory, and we may not want to completely rule them out by impose very strong conditions, and hence simple L^p bounds seem appropriate. Moreover, as we are dealing with compact manifolds, the L^1 norm is the weakest of all L^p norms and so the weakest possible L^p bound is to require $\varphi \in S \cap B_\eta(0)$.

At the end of the day, our hypotheses should be justified on physical grounds. Imposing that $\int_M \varphi \, dV_{g_0}$ is bounded above is more than appropriate, as we know from experience that if the extra compact dimensions exist, they have to be small. On the other hand, if $\int_M \varphi \, dV_{g_0} \to -\infty$ then such extra dimensions would collapse. It would be interesting to analyze what type of physics can emerge in this setting. For example, we illustrate in section 3.1 that, if the effective potential is not bounded from below then the singularities assume a rather specific form.

3. Proof of theorem 2.1

In this section we prove theorem 2.1. We use the letter C to denote several different constants that appear in the estimates. They will never depend on φ , u or i, where i indexes a sequence in S(see proof below). For $0 \leq \beta \leq 1$, the most natural choices are $\beta = 0, 1$, as we stated in the theorem. But since the proof works without modifications for values of β between zero and one, we will write $0 \leq \beta \leq 1$ explicitly in some passages of the proof in order to stress this fact.

Proof of theorem 2.1-(1): First we need to show that the map $\mathfrak{F}_{g_0,T_{st}^{g_0},\mathcal{F}}$ is well-defined. For a given $\varphi \in \mathcal{S}$, let w belong to the kernel of P_g , or equivalently to the kernel of M_{g_0} , where P_g and M_{g_0} have been defined in (2.6) and (2.15), respectively. By the Fredholm alternative (see appendix), w is then L^2 orthogonal to the image of P_g , and since -1 belongs to the image by hypothesis,

$$-\int_M w \, dV_g = -\int_M w e^{n\varphi} \, dV_{g_0} = 0.$$

Now, if u_1 and u_2 are two different solutions of (2.6), then $u_1 - u_2$ belongs to the kernel of P_q and

$$\int_{M} e^{n\varphi} u_1 \, dV_{g_0} = \int_{M} e^{n\varphi} u_2 \, dV_{g_0} + \int_{M} e^{n\varphi} (u_1 - u_2) \, dV_{g_0} = \int_{M} e^{n\varphi} u_2 \, dV_{g_0},$$

showing that the map is well defined.

Proof of theorem 2.1-(2), $0 \le \beta \le 1$: Here we prove part (2) of the theorem for β between zero and one. A separate proof for the case $\beta = \frac{n}{2}$ and when (2.14) holds will be provided for the reasons explained in section 2.

We start deriving an useful inequality. Notice that v > 0, and so is u. We are therefore allowed to divide equation (2.15) by u. Doing so and integrating by parts yields

(3.1)
$$\int_{M} \frac{|\nabla_{g_{0}}u|_{g_{0}}^{2}}{u^{2}} dV_{g_{0}} + \frac{1}{3}(n-1)(n-2) \int_{M} |\nabla_{g_{0}}\varphi|^{2} dV_{g_{0}} + (n-2) \int_{M} \frac{1}{u} \langle \nabla_{g_{0}}\varphi, \nabla_{g_{0}}u \rangle_{g_{0}} dV_{g_{0}} + \int_{M} \frac{e^{2\varphi}}{u} dV_{g_{0}} + \int_{M} \mathcal{F} dV_{g_{0}} - \frac{1}{6} \int_{M} e^{2(1-\beta)\varphi} T_{\mathrm{st}}^{g_{0}} dV_{g_{0}} = \frac{1}{3} R_{g_{0}} \operatorname{vol}_{g_{0}}(M).$$

The Cauchy inequality-with-epsilon (see appendix) gives

$$\int_M \frac{1}{u} \langle \nabla_{g_0} \varphi, \nabla_{g_0} u \rangle_{g_0} \, dV_{g_0} \ge -\frac{\varepsilon}{2} \int_M |\nabla_{g_0} \varphi|^2 \, dV_{g_0} - \frac{1}{2\varepsilon} \int_M \frac{|\nabla_{g_0} u|_{g_0}^2}{u^2} \, dV_{g_0}$$

 $\mathrm{so},$

(3.2)
$$\int_{M} \frac{|\nabla_{g_{0}}u|_{g_{0}}^{2}}{u^{2}} dV_{g_{0}} + \frac{1}{3}(n-1)(n-2) \int_{M} |\nabla_{g_{0}}\varphi|^{2} dV_{g_{0}} + (n-2) \int_{M} \frac{1}{u} \langle \nabla_{g_{0}}\varphi, \nabla_{g_{0}}u \rangle_{g_{0}} dV_{g_{0}} \\ \geq \left(1 - \frac{n-2}{2\varepsilon}\right) \int_{M} \frac{|\nabla_{g_{0}}u|_{g_{0}}^{2}}{u^{2}} dV_{g_{0}} + (n-2) \left(\frac{n-1}{3} - \frac{\varepsilon}{2}\right) \int_{M} |\nabla_{g_{0}}\varphi|^{2} dV_{g_{0}}.$$

Since $\frac{n-2}{2} < \frac{2(n-1)}{3}$ one can choose ε such that $\varepsilon > \frac{n-2}{2}$ and $\varepsilon < \frac{2(n-1)}{3}$. Hence (3.1) and (3.2) combine to give

(3.3)
$$\frac{1}{3}R_{g_0}\operatorname{vol}_{g_0}(M) \ge c_1 \int_M \frac{|\nabla_{g_0}u|_{g_0}^2}{u^2} dV_{g_0} + c_2 \int_M |\nabla_{g_0}\varphi|^2 dV_{g_0} + \int_M \frac{e^{2\varphi}}{u} dV_{g_0} + \int_M \mathcal{F} dV_{g_0} - \frac{1}{6} \int_M e^{2(1-\beta)\varphi} T_{\mathrm{st}}^{g_0} dV_{g_0},$$

where c_1 and c_2 are positive constants.

To prove (2.19), assume the result is not true. Then there exists a sequence $\varphi_i \in S_\eta$, with corresponding solutions u_i of (2.15), such that $\mathfrak{F}_{g_0,T^{g_0}_{\mathrm{st}}}(\varphi_i) \to -\infty$ as $i \to \infty$, and therefore

(3.4)
$$\int_M e^{n\varphi_i} u_i \, dV_{g_0} \to 0^+$$

We will suppress the subscript *i* for notational convenience, but all limits are to be understood as the limit when $i \to \infty$.

Start noticing that $\int_M e^{n\varphi} dV_{g_0} = A$ implies that there exists a constant C > 0 such that

$$\int_M e^{2(1-\beta)\varphi} \, dV_{g_0} \le C.$$

This is obvious when $\beta = 1$ and for $\beta < 1$, Hölder's inequality gives

$$\int_{M} e^{2(1-\beta)\varphi} dV_{g_{0}} \leq \left(\operatorname{vol}_{g_{0}}(M) \right)^{\frac{n}{n-2(1-\beta)}} \| e^{2(1-\beta)\varphi} \|_{L^{\frac{n}{2(1-\beta)}}(M,g_{0})} \\ = \left(\operatorname{vol}_{g_{0}}(M) \right)^{\frac{n}{n-2(1-\beta)}} \left(\int_{M} e^{n\varphi} dV_{g_{0}} \right)^{\frac{2(1-\beta)}{n}}.$$

Hence

$$\left| \int_{M} e^{2(1-\beta)\varphi} T_{\mathrm{st}}^{g_0} \, dV_{g_0} \right| \le \sup_{x \in M} |T_{\mathrm{st}}^{g_0}(x)| \int_{M} e^{2(1-\beta)\varphi} \, dV_{g_0} \le C \sup_{x \in M} |T_{\mathrm{st}}^{g_0}(x)|.$$

Combining this with the fact that the first, second and fourth integrals on the right hand side of (3.3) are non-negative, produces

$$C + \frac{1}{3}R_{g_0}\operatorname{vol}_{g_0}(M) \ge \frac{1}{6}\int_M e^{2(1-\beta)\varphi}T_{\mathrm{st}}^{g_0}\,dV_{g_0} + \frac{1}{3}R_{g_0}\operatorname{vol}_{g_0}(M) \ge \int_M \frac{e^{2\varphi}}{u}\,dV_{g_0}.$$

Or, in other words, $\int_M \frac{e^{2\varphi}}{u} dV_{g_0}$ is bounded independent of φ . Applying Cauchy-Schwartz gives

$$0 \leq \int_M e^{\frac{n+2}{2}\varphi} dV_{g_0} = \int_M e^{\frac{n}{2}\varphi} \sqrt{u} \frac{e^{\varphi}}{\sqrt{u}} dV_{g_0} \leq \left(\int_M e^{n\varphi} u \, dV_{g_0}\right)^{\frac{1}{2}} \left(\int_M \frac{e^{2\varphi}}{u} \, dV_{g_0}\right)^{\frac{1}{2}}.$$

As the last term is bounded, from (3.4) it follows that

$$\| e^{\varphi} \|_{L^{\frac{n+2}{2}}(M,g_0)}^{\frac{n+2}{2}} = \int_M e^{\frac{n+2}{2}\varphi} \to 0.$$

By the Hölder inequality it then follows that $|| e^{\varphi} ||_{L^s(M,g_0)} \to 0$ for any $1 \le s \le \frac{n+2}{2}$. Now we will bootstrap to obtain this result for s < n. For any p > 1, letting $\frac{1}{p} + \frac{1}{q} = 1$ as usual, we see that:

$$(3.5) \quad \int_{M} e^{s\varphi} \, dV_{g_0} = \int_{M} e^{(s-\frac{n}{p})\varphi} e^{\frac{n}{p}\varphi} \, dV_{g_0} \le \left(\int_{M} e^{q(s-\frac{n}{p})\varphi}\right)^{\frac{1}{q}} \left(\int_{M} e^{n\varphi}\right)^{\frac{1}{p}} = A^{\frac{1}{p}} \left(\int_{M} e^{\frac{p}{p-1}(s-\frac{n}{p})\varphi}\right)^{\frac{p-1}{p}}$$

Given $\frac{n+2}{2} < s < n$, we can always choose p > 1 such that

(3.6)
$$\frac{p}{p-1}(s-\frac{n}{p}) = \frac{n+2}{2},$$

so that the right hand side of (3.5) goes to zero. In fact, as p > 1, write $p = 1 + \delta$, $\delta > 0$, so that (3.6) reads

$$s = \frac{n+2}{2}\frac{\delta}{1+\delta} + \frac{n}{1+\delta}$$

Then s is a decreasing function of δ , satisfying $s \to n$ when $\delta \to 0$ and $s \to \frac{n+2}{2}$ as $\delta \to \infty$. Hence we have shown:

(3.7)
$$\int_{M} e^{s\varphi} \, dV_{g_0} \to 0 \quad \text{for any} \quad 1 \le s < n$$

Now multiply (2.15) by $\frac{e^{(n-2)\varphi}}{u}$, integrate, and integrate by parts the Laplacian terms. Then terms involving $\frac{1}{u} \langle \nabla_{g_0} \varphi, \nabla_{g_0} u \rangle_{g_0}$ cancel out and we obtain

(3.8)
$$\frac{1}{3}(n-1)(n-2)\int_{M}e^{(n-2)\varphi}|\nabla_{g_{0}}\varphi|^{2}\,dV_{g_{0}} = \int_{M}\frac{e^{n\varphi}}{u}\,dV_{g_{0}} + \int_{M}\frac{|\nabla_{g_{0}}u|^{2}}{u^{2}}e^{(n-2)\varphi}\,dV_{g_{0}} - \frac{1}{3}\int_{M}e^{(n-2)\varphi}R_{g_{0}}\,dV_{g_{0}} + \int_{M}\mathcal{F}\,dV_{g_{0}} - \frac{1}{6}\int_{M}e^{(n-2\beta)\varphi}T_{\mathrm{st}}^{g_{0}}\,dV_{g_{0}}.$$

If $\beta = 0$ then

$$\left|\int_{M} e^{(n-2\beta)\varphi} T_{\mathrm{st}}^{g_{0}} \, dV_{g_{0}}\right| = \left|\int_{M} e^{n\varphi} T_{\mathrm{st}}^{g_{0}} \, dV_{g_{0}}\right| \le \sup_{x \in M} |T_{\mathrm{st}}^{g_{0}}(x)| \int_{M} e^{n\varphi} \, dV_{g_{0}} = A \sup_{x \in M} |T_{\mathrm{st}}^{g_{0}}(x)|,$$

which is uniformly bounded independent of φ . If $\beta > 0$ (and less or equal than one, see the hypotheses of the theorem) then

$$\left| \int_{M} e^{(n-2\beta)\varphi} T_{\mathrm{st}}^{g_{0}} \, dV_{g_{0}} \right| = \sup_{x \in M} |T_{\mathrm{st}}^{g_{0}}(x)| \int_{M} e^{(n-2\beta)\varphi} \, dV_{g_{0}}$$

which goes to zero by (3.7). In any situation, there exists a constant C > 0, independent of φ such that

(3.9)
$$-\frac{1}{6} \int_{M} e^{(n-2\beta)\varphi} T_{\mathrm{st}}^{g_0} \, dV_{g_0} \ge -C \sup_{x \in M} |T_{\mathrm{st}}^{g_0}(x)|$$

Evoking (3.7) again gives

(3.10)
$$\frac{1}{3} \int_{M} e^{(n-2)\varphi} R_{g_0} \, dV_{g_0} \to 0$$

Using (3.9), (3.10) with

$$\int_M \frac{|\nabla_{g_0} u|^2}{u^2} \, dV_{g_0} \ge 0 \quad \text{and} \quad \int_M \mathcal{F} \, dV_{g_0} \ge 0$$

into (3.8) implies

(3.11)
$$\frac{1}{3}(n-1)(n-2)\int_{M} e^{(n-2)\varphi} |\nabla_{g_0}\varphi|^2 \, dV_{g_0} \ge \int_{M} \frac{e^{n\varphi}}{u} \, dV_{g_0} - C \sup_{x \in M} |T_{\mathrm{st}}^{g_0}(x)| + o(1),$$

where o(1) denotes as usual a term that goes to zero. Since u > 0 and $f(x) = \frac{1}{x}$ is convex for x > 0, Jensen's inequality (see appendix) with the measure $dV_g = e^{n\varphi} dV_{g_0}$ gives:

(3.12)
$$\int_{M} \frac{e^{n\varphi}}{u} \, dV_{g_0} = \operatorname{vol}_g(M) \oint_{M} \frac{1}{u} \, dV_g \ge A^2 \frac{1}{\int_{M} u \, dV_g} = \frac{A^2}{\int_{M} e^{n\varphi} u \, dV_{g_0}}$$

where $f_M = \frac{1}{\text{vol}(M)} \int_M$. Therefore $\int_M \frac{e^{n\varphi}}{u} dV_{g_0}$ goes to infinity, and so by (3.11) we get

(3.13)
$$\int_{M} e^{(n-2)\varphi} |\nabla_{g_0}\varphi|^2 \, dV_{g_0} \to \infty.$$

Now integrate (2.9) with respect to $dV_g = e^{n\varphi} dV_{g_0}$ and integrate by parts the Laplacian term to find

$$\int_{M} R_g \, dV_g = (n-1)(n-2) \int_{M} e^{(n-2)\varphi} |\nabla_{g_0}\varphi|^2 \, dV_{g_0} + R_{g_0} \int_{M} e^{(n-2)\varphi} \, dV_{g_0}.$$

The second integral approaches zero by (3.7), and since $\int_M R_g dV_g \leq \eta$ by hypothesis, we obtain a contradiction with (3.13). This proves (2.19) when $0 \leq \beta \leq 1$.

Proof of theorem 2.1-(2), $\beta = \frac{n}{2}$ or (2.14) holds: It is enough to assume (2.14), since this is satisfied when $\beta = \frac{n}{2}$ because

$$\int_M T_{\mathrm{st}}^g \, dV_g = \int_M e^{-n\varphi} T_{\mathrm{st}}^{g_0} e^{n\varphi} \, dV_g = \int_M T_{\mathrm{st}}^{g_0} \, dV_{g_0}.$$

As before, assume the result is not true so that

(3.14)
$$\int_M e^{n\varphi} u \, dV_{g_0} = \int_M u \, dV_g \to 0$$

for some sequence of functions u and φ (again we are omitting the index i). Here it will be more convenient to work with equation (2.6). Divide (2.6) by u, integrate with respect to dV_g and integrate by parts the Laplacian term to get

(3.15)
$$\int_{M} \frac{|\nabla_{g}u|^{2}}{u^{2}} dV_{g} + \int_{M} \left(-\frac{1}{3}R_{g} + \frac{1}{6}F_{g} - \frac{1}{6}T_{st}^{g}\right) dV_{g} = -\int_{M} \frac{1}{u} dV_{g}$$

By (2.14), the integral

$$\int_M T_{\rm st}^g \, dV_g$$

is bounded by some constant C > 0 depending only on the fixed quantities $T_{\rm st}^{g_0}$ and g_0 . Hence, dropping out the non-negative terms,

(3.16)
$$\int_M \frac{1}{u} dV_g \le \int_M \frac{1}{3} R_g \, dV_g + \int_M T_{\rm st}^g \, dV_g \le C,$$

where the hypothesis that the integral of the scalar curvature is bounded has been used. As before, applying Jensen's inequality for the function $f(x) = \frac{1}{x}$, x > 0, and the measure dV_g leads to

$$\frac{1}{f_M \, u \, dV_g} \le \int_M \frac{1}{u} dV_g,$$

or equivalently

$$\frac{1}{\int_M u \, dV_g} \le \frac{1}{A^2} \int_M \frac{1}{u} dV_g,$$

where we used $\int_M dV_g = \int_M e^{n\varphi} dV_{g_0} = A$. In light of (3.14), the right hand side of the above inequality goes to infinity, contradicting (3.16). This proves (2.19) when $\beta = \frac{n}{2}$ or (2.14) holds.

The statement (2.20) now follows from the Gauss-Bonnet formula as $\int_M R_g dV_g$ is a topological invariant, and hence bounded.

Proof of theorem 2.1-(3): Use again (3.15). It then follows from the hypothesis that

$$\int_M \frac{1}{u} \, dV_g \le C_g$$

which again gives a contradiction after an application of the Jensen's inequality.

Proof of proposition 2.2: If the result is not true then from the proof of theorem 2.1 we have $\int_M e^{s\varphi} dV_{g_0} \to 0$ for $1 \leq s < n$. Since the exponential is a convex function, we can use Jensen's inequality with the measure dV_{g_0} to obtain:

$$\int_{M} e^{s\varphi} \, dV_{g_0} = \operatorname{vol}_{g_0}(M) \oint_{M} e^{s\varphi} \, dV_{g_0} \ge \operatorname{vol}_{g_0}(M) e^{s \frac{1}{\operatorname{vol}_{g_0}(M)} \int_{M} \varphi \, dV_{g_0}},$$

where $\oint_{M} = \frac{1}{\operatorname{vol}_{g_0}} \int_{M}$. So $\int_{M} e^{s\varphi} \to 0$ implies $\int \varphi \, dV_{g_0} \to -\infty$, and therefore
 $\infty \leftarrow \left| \int \varphi \, dV_{g_0} \right| \le \int |\varphi| \, dV_{g_0} = \| \varphi \|_{L^1(M,g_0)},$

contradicting $\varphi \in \mathcal{S} \cap B_n(0)$.

3.1. **Remarks.** The inequality Eq. (3.3) is also interesting as a limit on how negatively curved the compactification manifold M can be. In [15] it was pointed out that M cannot have negative curvature unless $T_{st} > 0$, but this was not quantified. From equation (3.3) one has (see also proposition 4.2)

$$R_{g_0} \ge -\frac{\int_M e^{2(1-\beta)\varphi} T_{\rm st}^{g_0} \, dV_{g_0}}{2 \operatorname{vol}_{g_0}(M)}$$

Thus the negative curvature can be no greater than string scale, and much less if the orientifolds live on submanifolds.

We mentioned in section 2 that if the effective potential is not bounded from below, then the singularities that may form are generally of a very particular type. To see this, remember that we showed that in this setting we have

(3.17)
$$\int_M e^{s\varphi_i} \, dV_{g_0} \to 0,$$

when $i \to \infty$, for any $1 \le s < n$. Since $\int_M e^{s\varphi_i} dV_{g_0} = A$, for all *i*, a simple application of the interpolation inequality (see appendix) yields

(3.18)
$$\int_M e^{s\varphi_i} \, dV_{g_0} \to \infty$$

for any s > n. On the other hand, (3.17) gives that, up to a subsequence, $e_i^{\varphi} \to 0$ almost everywhere, and by Egoroff's theorem (see e.g. [34]) the convergence is uniform outside a set of measure δ , where $\delta > 0$ is as small as desired. This implies that $e^{n\varphi_i}$ has a behavior very much like a (sum of) Dirac delta(s): it converges to zero in most of M, blows-up in some localized subsets and has constant integral (since $\int_M e^{n\varphi_i} dV_{g_0} = A$).

4. EXISTENCE OF CRITICAL POINTS, TOPOLOGICAL CONSIDERATIONS AND SOME A PRIORI ESTIMATES

In this section we discuss the existence of critical points and its interplay with the topology/geometry of the manifold (M, g_0) . First we prove existence theorems for equation (2.15) which also allow us to provide useful bounds for the critical points and the map $\mathfrak{F}_{g_0,T^{g_0}_{st},\mathcal{F}}$. We then discuss the general solvability of equation (2.1), including non-positive values of α .

Recall from $\S2$ that our problem is to solve equation (2.1),

$$P_g v = \Delta_g v + (-\frac{1}{3}R_g + \frac{1}{6}F_g - \frac{1}{6}T_{\rm st}^g)v = \frac{1}{6}\alpha$$

for some $\alpha \in \mathbb{R}$ and an everywhere positive function v satisfying equation (2.2),

$$\int_M v \, dV_g = \frac{1}{G_N}.$$

Of course, if we can solve (2.1) for some α , then by simultaneously rescaling α and v we can solve (2.2). Thus, for fixed data (g, F, T_{st}) , we can restrict attention to the cases $\alpha = \pm 1$ or 0, as in equation (2.6). On the other hand, if we are considering a family of solutions in which α changes sign, we should keep α general.

The solvability of (2.1) is governed by the associated eigenvalue problem

$$P_q v = \lambda v$$

We recall its basic features, granting that the coefficient functions are smooth, Δ_g is negative definite and that P_g is self-adjoint (see e.g. [16, 19]):

- The spectrum Spec $P_g \subset \mathbb{R}$ is countable and has no accumulation points.
- It is bounded above, Spec $P_g \subset (-\infty, \lambda_0]$.
- The eigenvalue λ_0 is isolated, meaning Spec $P_q = (-\infty, \lambda_1] \cup \{\lambda_0\}$ with $\lambda_1 < \lambda_0$.
- The eigenspace to λ_0 is one dimensional and spanned by an everywhere positive eigenfunction.
- The inhomogeneous problem $(P_g \mu)u = f$ has a unique solution for each $f \in C^{\infty}(M)$ if and only if $\mu \notin \operatorname{Spec} P_q$. Such a solution u is smooth.

Thus, if $\lambda_0 < 0$, P_g will have no kernel, so that existence and uniqueness of u is evident. On the other hand, if $\lambda_0 \ge 0$ or if we cannot show $\lambda_0 < 0$ a priori, the discussion will be more complicated. In this case we will try to argue by continuation from the easier $\lambda_0 < 0$ case.

Let us first consider cases where we can state conditions which guarantee $\lambda_0 < 0$, which we do in propositions 4.1 and 4.2. The idea behind their proofs is to guarantee the existence of a maximum principle, which can then be combined with the Fredholm alternative to reduce the problem of existence of solutions to that of uniqueness⁶ (see the appendix for details).

For simplicity, let us assume that β is an integer and that the gauge fields $F^{(p)}$ are labeled according to their degree, which we still denote by p, so that $p = 1, \ldots, n$. The proofs can easily be extended to other cases.

Proposition 4.1. Assume the same hypotheses of theorem 2.1, and suppose further that (M, g_0) has positive Yamabe invariant. Define

$$K_{g_0,n}(\varepsilon) = \frac{R_{g_0} - 3\varepsilon}{3n^n}.$$

Consider the conditions:

- (i) $|| T_{st}^{g_0} ||_{C^0(M)} < K_{g_0,n}(\varepsilon).$
- $\begin{array}{c} (ii) \parallel |F^{(p)}|_{g_0}^2 \parallel_{C^0(M)} < K_{g_0,n}(\varepsilon), \ p = 1, \dots, n. \\ (iii) \parallel \Delta_{g_0} \varphi \parallel_{C^0(M)} < K_{g_0,n}(\varepsilon). \end{array}$

- (iv) $\| |\nabla_{g_0} \varphi|_{g_0}^2 \|_{C^0(M)} \leq K_{g_0,n}(\varepsilon).$ (v) $\frac{1}{n} \leq e^{2\varphi(x)} \leq n^n$ for all $x \in M$.

Given $0 \le \varepsilon < \frac{1}{3}R_{g_0}$, assume conditions (i)-(v). Then equation (2.15) has a unique solution u. Such solution is smooth and positive. Moreover, if $\varepsilon > 0$, this solution obeys the estimate

(4.1)
$$\| u \|_{C^0(M)} \leq \frac{n^n}{\varepsilon}$$

Also, if $\varepsilon = 0$ then the map $\mathfrak{F}_{g_0,T_{st}^{g_0},\mathcal{F}}$ satisfies

(4.2)
$$\mathfrak{F}_{g_0,T^{g_0}_{st},\mathcal{F}}(\varphi) > -\frac{3n(n+3)R_{g_0}}{2AG_N^2}$$

where $K_{a_0,n} = K_{a_0,n}(0)$.

In order to appreciate the relevance of proposition 4.1, let us consider one of the cases of primary interest, namely, D = 10, so that n = 6, and $K_{g_0,6} \sim 10^{-5} R_{g_0}$. Here all of R_{g_0} , $T_{st}^{g_0}$ and \mathcal{F} have units of $1/length^2$, or mass² in Planck units $\hbar = c = 1$. Assume for concreteness that M is roughly isotropic in the sense that $\operatorname{vol}_{g_0} M \sim (R_{g_0})^{-6/2}$. Then proposition 4.1 guarantees existence and uniqueness of a critical point for any value of of $T_{\rm st}^{g_0}$ and \mathcal{F} up to the order of $10^{-5} (M_P^{(10)})^2$.

The known applicability of this proposition is to compactifications on manifolds with Einstein metrics of positive scalar curvature, such as the sphere. These lead to $\alpha < 0$ and compactifications to anti-de Sitter space-time. However it is interesting to note that when n = 6 and M is a Calabi-Yau manifold, it is known that its topological Yamabe invariant is positive (proposition 4.3 in [25]). Hence, in such cases, there exists at least one conformal class whose Yamabe invariant is positive and therefore there are solutions to equation (2.15) with positive scalar curvature, very different from the Ricci flat metrics usually considered. A physics discussion of this is given in [20].

 $^{^{6}}$ It is only the sign of the lower order term which plays a crucial role in this procedure, but rather than simply assuming such a sign condition, we provide separate bounds in terms of the physical quantities composing the lower order coefficient, namely, the scalar curvature, the string term and the gauge fields. Since it is only the combination of these terms which is relevant for the maximum principle argument, the choice of such bounds involves a great deal of arbitrariness, but the reader can easily adapt the proof to other situations. We also seek to write conditions in a very simple fashion (we provide a bound essentially in terms of the dimension), so that one can hope to actually verify them in concrete situations.

One interesting consequence of (4.1) is to give a control of critical points in terms of the (constant) scalar curvature R_{g_0} . In fact, choosing $\varepsilon = \frac{1}{4}R_{g_0} < \frac{1}{3}R_{g_0}$, (4.1) gives

(4.3)
$$\| u \|_{C^0(M)} \le \frac{4n^n}{R_{g_0}},$$

and therefore u has to be very small if the scalar curvature is very large. This behavior of u has already been identified in [14] by an argument based on the AdS/CFT correspondence, and (4.3) provides a further refinement of that prediction.

Proof of proposition 4.1: Write the operator M_{g_0} as

(4.4)
$$M_{g_0}u = \Delta_{g_0}u + (n-2)\langle \nabla_{g_0}\varphi, \nabla_{g_0}u\rangle_{g_0} - Vu$$

where $V = -\frac{2}{3}(n-1)\Delta_{g_0}\varphi - \frac{1}{3}(n-1)(n-2)|\nabla_{g_0}\varphi|_{g_0}^2 + \frac{1}{3}R_{g_0} - \mathcal{F} + \frac{1}{6}e^{2(1-\beta)\varphi}T_{\mathrm{st}}^{g_0}$. We will show that conditions (i) - (v) imply

(4.5)
$$\inf_{M} V(x) > \varepsilon.$$

We have

(4.6)
$$V > \frac{1}{3}R_{g_0} - \frac{2}{3}(n-1)K_{g_0,n}(\varepsilon) - \frac{1}{3}(n-1)(n-2)K_{g_0,n}(\varepsilon) - \frac{1}{6}\sum_{p=1}^{L}e^{2(1-p)\varphi}K_{g_0,n}(\varepsilon) - \frac{1}{6}e^{2(1-\beta)\varphi}K_{g_0,n}(\varepsilon).$$

Since we are assuming β to be an integer, we can write the last two terms into a single sum which is greater than or equal to

(4.7)

$$-\frac{2}{6}K_{g_0,n}(\varepsilon)\sum_{\ell=0}^{n} \left(e^{2\varphi}\right)^{1-\ell} \ge -\frac{1}{3}K_{g_0,n}(\varepsilon)\left(1+e^{2\varphi}+\sum_{\ell=2}^{n} \left(\frac{1}{e^{2\varphi}}\right)^{\ell-1}\right) \\
\ge -\frac{1}{3}K_{g_0,n}(\varepsilon)\left(1+n^n+(n-2)n^{n-1}\right) \ge -\frac{2}{3}n^n K_{g_0,n}(\varepsilon),$$

where in the next to the last step we used (v). Using (4.7) in (4.6) gives

$$V > \frac{1}{3}R_{g_0} - \frac{1}{3}n(n-1)K_{g_0,n}(\varepsilon) - \frac{1}{3}n^n K_{g_0,n}(\varepsilon)$$

$$\geq \frac{1}{3}R_{g_0} - \frac{1}{3}n^n K_{g_0,n}(\varepsilon) - \frac{2}{3}n^n K_{g_0,n}(\varepsilon) = \frac{1}{3}R_{g_0} - n^n K_{g_0,n}(\varepsilon).$$

The right hand side is equal to ε by the definition of $K_{g_0,n}(\varepsilon)$, and so (4.5) follows. In particular V > 0, including when $\varepsilon = 0$ (we get a strict inequality by the compactness of M), and it then follows that M_{g_0} satisfies the maximum principle (see appendix), and therefore M_{g_0} has trivial kernel. We give the argument for completeness. If $M_{g_0}w = 0$ then, from $M_{g_0}w \ge 0$ we get by the maximum principle that w cannot have a non-negative maximum, hence $w \le 0$. Analogously $M_{g_0}w \le 0$ implies that w cannot have a non-positive minimum, hence $w \ge 0$ and therefore $w \equiv 0$. Now by the Fredholm alternative, (4.4) has a unique solution u, and this solution is is smooth. To see that u > 0, we again evoke the maximum principle. We have $M_{g_0}u = -e^{2\varphi} \le 0$. Hence u cannot have a non-positive minimum.

To prove (4.1), define $\hat{u} = u - \frac{n^n}{\varepsilon}$. Using (4.4),

$$M_{g_0}\widehat{u} = \Delta_{g_0}\widehat{u} + (n-2)\langle \nabla_{g_0}\varphi, \nabla_{g_0}\widehat{u}\rangle_{g_0} - V\widehat{u} = -e^{2\varphi} + \frac{n^n}{\varepsilon}V > -e^{2\varphi} + n^n \ge 0,$$

where (4.5) and (v) have been used. Evoking the maximum principle once more we obtain $\hat{u} \leq 0$, and (4.1) follows.

Next we show (4.2). Let $x_0 \in M$ be a point where u attains its minimum. Then $\Delta_{g_0} u(x_0) \ge 0$ and $\nabla_{g_0} u(x_0) = 0$. Therefore

$$-e^{2\varphi(x_0)} = M_{g_0}u(x_0) \ge -V(x_0)u(x_0),$$

and so

$$u(x_0) \ge \frac{e^{2\varphi(x_0)}}{V(x_0)} \ge \frac{1}{n \sup_M V}$$

But

$$V \leq |V| \leq \frac{2}{3}(n-1)|\Delta_{g_0}\varphi| + \frac{1}{3}(n-1)(n-2)|\nabla_{g_0}\varphi|_{g_0}^2 + \frac{1}{3}R_{g_0} + \mathcal{F} + \frac{1}{6}e^{2(1-\beta)\varphi}|T_{\rm st}^{g_0}| \\ < \frac{1}{3}n(n-1)K_{g_0,n} + \frac{1}{3}R_{g_0} + \frac{2}{6}K_{g_0,n}\sum_{\ell=0}^n \left(e^{2\varphi}\right)^{1-\ell} \\ \leq \frac{1}{3}n(n-1)K_{g_0,n} + \frac{1}{3}R_{g_0} + \frac{1}{3}K_{g_0,n}(n+1)n^n \leq (n+3)R_{g_0}$$

where the definition of K_{n,g_0} has been used. Hence

$$u(x_0) \ge \frac{1}{n(n+3)R_{g_0}}$$

Now it follows that

$$\mathfrak{F}_{g_0,T_{\rm st}^{g_0},\mathcal{F}}(\varphi) = -\frac{6}{4G_N^2 \int_M e^{n\varphi} u \, dV_{g_0}} \ge -\frac{6}{4G_N^2 \, \inf_M u \, \int_M e^{n\varphi} \, dV_{g_0}} \ge -\frac{3n(n+3)R_{g_0}}{2AG_N^2},$$

where we used that $\int_M e^{n\varphi} \, dV_{g_0} = A.$

4.1. Non-positive Yamabe invariant.

Proposition 4.2. Assume the same hypotheses of theorem 2.1, and suppose further that (M, g_0) has non-positive Yamabe invariant. Let $H = \frac{1}{6n^{|1-\beta|}}(-n^n \sum |F^{(p)}|_{g_0}^2 + T_{st}^{g_0})$. Given $\varepsilon > 0$, and $\Gamma > 1$, define

$$K_{g_0,n}(\varepsilon,\Gamma) = \frac{n^2\Gamma}{3} + \frac{|R_{g_0}|}{3} + \varepsilon.$$

Assume that

$$\| \Delta_{g_0} \varphi \|_{C^0(M)} < \Gamma$$
$$\| |\nabla_{g_0} \varphi|_{g_0}^2 \|_{C_0(M)} < \Gamma$$
$$\frac{1}{n} \le e^{2\varphi} \le n$$
$$H > K_{g_0,n}(\varepsilon, \Gamma)$$

Then, there exists a unique, smooth and, positive solution to (2.15). This solution satisfies

(4.8)
$$\| u \|_{C^0(M)} \leq \frac{n}{\epsilon}.$$

16

Moreover, the map $\mathfrak{F}_{q_0,T^{g_0}_{st},\mathcal{F}}$ satisfies,

(4.9)
$$\mathfrak{F}_{g_0,T_{st}^{g_0},\mathcal{F}} > -\frac{n^{1+|1-\beta|} \|T_{st}^{g_0} - \frac{1}{n^n} \sum_{p=1}^n |F^{(p)}|^2 \|_{C^0(M)}}{2G_N^2 A}$$

It is known that there are no topological obstructions to negative scalar curvature [5]. In other words, on any compact manifold we can find a metric g such that the total scalar curvature functional is negative, i.e., $\int_M R_g dV_g < 0$, so (M,g) will have negative Yamabe invariant. Thus we have conditions which guarantee a solution in this case; however they may require large $T_{\rm st}^{g_0}$ which is somewhat unphysical. Nevertheless this could be useful as a starting point for the analysis of families of solutions below.

For the situation of zero Yamabe invariant, we notice that once more in the case of primary interest, namely, when M is a Calabi-Yau, we can find a metric fulfilling the hypotheses of proposition 4.2. Simply choose a Ricci-flat metric on M. Therefore, as in proposition 4.1, our hypotheses in proposition 4.2 are not vacuous.

Proof of proposition 4.2: The proof is similar to proposition 4.1, so we will only sketch the arguments. Just as in the proof of proposition 4.1, we may easily show that the hypotheses imply (4.5). Then, the same maximum principle type of argument as before shows there exists a unique, smooth and positive solution to (2.15).

To show (4.8), if x_0 is a point of maximum of u, then, $\nabla_{g_0} u(x_0) = 0$ and $\Delta_{g_0} u(x_0) \leq 0$. Hence,

$$V(x_0)u(x_0) \le e^{2\varphi(x_0)}$$

So, $u \leq \frac{n}{\epsilon}$.

Finally, we prove (4.9). At a point of minimum y_0 of u, $\nabla_{g_0} u(y_0) = 0$ and $\Delta_{g_0} u(y_0) \ge 0$. Hence,

$$V(y_0)u(y_0) \ge e^{2\varphi(y_0)}$$

But,

$$V < \frac{1}{3} |R_{g_0}| + \frac{1}{3} n(n-1)\Gamma + \sup \frac{e^{2(1-\beta)\phi}}{6} |T_{st}^{g_0} - \sum_{p=1}^n |F^{(p)}|^2 e^{2(\beta-p)\phi}|$$

$$< \frac{n^{|1-\beta|}}{3} ||T_{st}^{g_0} - \frac{1}{n^n} \sum_{p=1}^n |F^{(p)}|^2 ||_{C^0(M)},$$

and the result follows.

4.2. Families of solutions. The basic point we want to establish is

Conjecture 4.3. Suppose we have fixed $(M, g_0, T_{st}^{g_0})$ and we are given data (φ, \mathcal{F}) depending smoothly on a real parameter t, such that for t = 0 we are in one of the situations governed by the preceding propositions and thus $\alpha|_{t=0} < 0$. Then, there will be a region $t \in [0, t_1]$ for some $t_1 > 0$ in which equation (2.1) has a unique solution satisfying (2.2), and in some cases $\alpha|_{t=t_1} > 0$.

In particular, we might add a constant term Ct to \mathcal{F} , causing the eigenvalues of P_g to shift as $\lambda \to \lambda + Ct$; then for a given t_1 and sufficiently large C we will have $\alpha|_{t=t_1} > 0$. Physically, this connects a compactification to anti-de Sitter space-time to another compactification to de Sitter space-time, and is referred to as "uplifting" [13, 21].

DISCONZI, DOUGLAS, AND PINGALI

Let us give some intuition of why one expects the above conjecture, or some variation of it, to be true. Begin by assuming that the corresponding solutions $v = v_{\alpha}$ vary in a well behaved manner (say, they depend continuously or smoothly on the data). What can be said about the solutions v_0 corresponding to $\alpha = 0$?

The (unique) solution of (2.1) can be written as an eigenfunction expansion of the form

(4.10)
$$v = \frac{\alpha}{6} \sum_{i} \frac{1}{\lambda_i} \left(\int_M \psi_i \, dV_g \right) \psi_i,$$

and the effective potential is then given by

(4.11)
$$\frac{1}{\mathcal{V}} = \frac{2}{3} G_N^2 \sum_i \frac{1}{\lambda_i} \left(\int_M \psi_i \, dV_g \right)^2.$$

Equality (4.11) indicates the presence of a resonant mode when \mathcal{V} approaches zero. In other words, when $\alpha \to 0$, one has $\lambda_s \to 0$ for one of the eigenvalues λ_s in $\{\lambda_i\}$; notice that by (4.10), if $\alpha \to 0$ then having $\lambda_s \to 0$ is the only way we can have a non-trivial v (at least under the present assumption that v will vary smoothly with the data of the problem). In this case, the main contribution to the warp factor will come from the corresponding eigenfunctions ψ_s , as can be seen from (4.10). It is also worth noticing that if $\lambda_s = 0$ then ψ_s will solve

$$P_q\psi_s = 0$$

which is exactly (2.1) with $\alpha = 0$.

In order to be an acceptable solution, one needs v to be positive. If ψ_s is the ground state, i.e., $\lambda_s = 0$ is the principal eigenvalue, then the corresponding eigenspace is one-dimensional and we can assume $\psi_s > 0$, and hence $v = \psi$. But if the resonant state is an excited one, then there might be more than one corresponding ψ_s , and generally these will take positive and negative values. Since the dominant contribution to v will be from the ψ_s 's, generally v will be zero or negative somewhere on M. Hence we seek to investigate whether this situation can be avoided, at least in some generic sense.

Suppose \mathcal{V} crosses from negative to positive values, and consider \mathcal{V} right before it becomes zero, so $\mathcal{V} \approx 0$ but $\mathcal{V} < 0$. Then, up to a positive constant, (4.11) becomes

$$\frac{1}{\mathcal{V}} \approx \frac{1}{\lambda_s},$$

where λ_s corresponds to the energy level that becomes resonant when $\mathcal{V} = 0$. The condition $\mathcal{V} < 0$ gives $\lambda_s < 0$, and therefore (4.10) implies

(4.12)
$$v \approx \frac{\alpha}{\lambda_s} \sum_{j=1}^k \psi_{s_k},$$

where ψs_k are the eigenfunctions associated to λ_s ; notice that $\frac{\alpha}{\lambda_s} > 0$ since $\mathcal{V} < 0 \Leftrightarrow \alpha < 0$.

Now recall that we are assuming that for $\alpha < 0$, or equivalently $\mathcal{V} < 0$, (2.1) has a (unique) solution v, which satisfies v > 0. As we pointed out earlier, typically one expects the eigenfunctions associated to excited states to take positive and negative values, and therefore the positivity of v along with (4.12) suggest that we are in the case in which λ_s is the principal eigenvalue, so that the sum on the right hand side of (4.12) contains only one term ψ_s and it satisfies $\psi_s > 0$. Therefore, under our assumptions, we expect that when $\mathcal{V} = 0$, the resonance $\lambda_s = 0$ will be the principal

eigenvalue, and in this case the warp factor will be $v = \psi_s > 0$ (up to multiplication to positive constants).

Summarizing, the above arguments suggest that typically one can continue \mathcal{V} through zero. Other arguments in this direction have already been given in [14]. Moreover, the study of continuous dependence of eigenvalues of the Laplacian [6] and analytic properties of the resolvent [27] suggest that a picture where some sort of smooth or continuous dependence, as assumed above, is likely to hold.

Finally, it is worth pointing out that if we rewrite all quantities in terms of a fixed constant scalar curvature background metric g_0 , then our basic equation becomes

$$L_{g_0}u - \frac{1}{3}R_{g_0}u = \frac{1}{6}e^{2\varphi}\alpha$$

and the study of eigenvalues previously discussed can be rephrased into the more geometric question of whether a multiple of the scalar curvature, $\frac{1}{3}R_{g_0}$, is an eigenvalue of the linear operator L_{g_0} .

5. The case $d \neq 4$

Although d = 4 for compactifications of string theory which could describe fundamental physics, there is also a good deal of physics work on $d \neq 4$. Such compactifications can be much simpler than d = 4, for example the compactification of the IIb superstring on S^5 leading to anti-de Sitter space-time [1]. A compactification with d > 4 can also be further compactified to get d = 4.

The reasoning of section 2, when applied to the general d case yields the following equation for the critical points of the action (see [14] and equation (2.33) therein for details)

(5.1)
$$\Delta_g v - \frac{1}{2(d-1)} \left(\frac{d}{2} R_g + T^{(d)} \right) v = \frac{(d-2)\alpha}{4(d-1)} v^{1-\frac{4}{d}}, \ v > 0,$$

where $T^{(d)} = -\frac{d}{2}F_g + T_{st}^g = -\frac{d}{2}\sum |F_g^{(p)}|^2 + T_{st}^g$ and α is the Lagrange multiplier in (1.1). One sees that the case d < 4 is qualitatively different because of the negative exponent in the source term, and we will not treat it here.⁷

Recalling that the role of α is to impose the constraint

(5.2)
$$\int_{M} v^{2-\frac{4}{d}} dV_g = \frac{1}{G_N},$$

we see that it suffices then to find v > 0 solving the problem

$$\Delta_g v + f_g v = K v^{1 - \frac{4}{d}},$$

where f_g is a given function, possibly depending on g, and K is a non-zero constant. In fact, if v solves (5.3) with $f = -\frac{1}{2(d-1)}(\frac{d}{2}R_g + T^{(d)})$, then $\tilde{v} = av$ satisfies (5.2) and solves the equation with K replaced by $Ka^{\frac{4}{d}}$, where

$$a = \left(\frac{1}{G_N \int_M v^{2-\frac{4}{d}} \, dV_g}\right)^{\frac{d}{2d-4}}$$

The important aspect of K is that it has the same sign of α . The case K = 0, or $\alpha = 0$, will be excluded here because in this case the equation is linear.

⁷ The case d = 2 is substantially different as one should use different constraints than those of Conjecture 1.1, *i.e.* one should not go to Einstein frame.

Let $L_g = \Delta_g + f_g$. Whenever a solution v > 0 of (5.3) exists, the sign of K is opposite to that of the lowest eigenvalue of L_g : let λ_1 be the lowest eigenvalue with corresponding eigenfunction $\phi_1 > 0$, and take the L^2 inner product to find, after integration by parts,

$$\int_{M} \phi_1 L_g v \, dV_g = \int_{M} v L_g \phi_1 \, dV_g = -\lambda_1 \int_{M} \phi_1 v \, dV_g = K \int_{M} \phi_1 v^{1-\frac{4}{d}} \, dV_g.$$

Hence, the sign of K, and therefore that of α , can be determined by studying the eigenvalue problem for the linear portion of the equation.⁸ Not surprisingly, then, different behaviors for solutions can be expected according to the sign of the first eigenvalue of L_g . With that in mind, we now turn to some general existence results.

Proposition 5.1. Assume f_g is smooth, K > 0 and d > 4. Then problem (5.3) has a smooth non-negative solution v satisfying

(5.4)
$$\|v\|_{C^0(M)} \le \left(\frac{K}{\|f_g\|_{C^0(M)}}\right)^{\frac{a}{4}}$$

Proof. Equation (5.3) can be solved by a standard sub- and super-solutions method (see proposition A.5). It is enough to find two functions $v_{-} < v_{+}$ such that

$$L_g v_- - K v_-^{1-\frac{4}{d}} \ge 0,$$

$$L_g v_+ - K v_+^{1-\frac{4}{d}} \le 0.$$

For v_{-} we can take $v_{-} \equiv 0$. If we choose v_{+} to be a positive constant, then

$$L_g v_+ - K v_+^{1-\frac{4}{d}} = f_g v_+ - K v_+^{1-\frac{4}{d}} \le \| f_g \|_{C^0(M)} v_+ - K v_+^{1-\frac{4}{d}}.$$

Setting the right hand side of the above expression to be less than or equal to zero implies

$$v_+ \le \left(\frac{K}{\parallel f_g \parallel_{C^0(M)}}\right)^{\frac{a}{4}}$$

Therefore one can set

$$v_+ \equiv \left(\frac{K}{\parallel f_g \parallel_{C^0(M)}}\right)^{\frac{a}{4}}.$$

Then, proposition A.5 guarantees the existence of a solution to (5.3) satisfying $v_{-} \equiv 0 \leq v \leq v_{+}$. In particular the bound (5.4) holds. Smoothness then follows from elliptic regularity (proposition A.3).

Remark 5.2. Without further information on f_g , we cannot guarantee that the solution v found in 5.1 is strictly positive. For example, if f_g vanishes identically then $v \equiv 0$ is the only non-negative solution to (5.3).

⁸In particular, if $T^{(d)}$ is somehow tuned so that $f_g = -\frac{n-2}{4(n-1)}R_g$, then L_g becomes the conformal Laplacian of the manifold (M, g), and therefore the sign of λ_1 , and hence of K, is an invariant of the conformal class [26, 4, 40].

For K < 0, we generally expect the negative part of f_g to dominate. Indeed, since if $v \ge 0$ then, integrating (5.3), gives

(5.5)
$$\int_M f_g v \, dV_g = K \int_M v^{1-\frac{4}{d}} \, dV_g \le 0$$

Furthermore, if $f_g < 0$ and v > 0 then, the integral on the left hand side of (5.5) is negative, and one sees that a necessary condition for the existence of non-negative solutions to (5.3) is K < 0. This motivates the following:

Proposition 5.3. Assume $f_g < 0$ is smooth, K < 0 and d > 4. Then problem (5.3) has a smooth positive solution v satisfying

(5.6)
$$\|v\|_{C^0(M)} \le \left(\frac{|K|}{\min_M |f_g|}\right)^{\frac{d}{4}}$$

Remark 5.4. Recall that, in the case of interest $f = -\frac{1}{2(d-1)}(\frac{d}{2}R_g + T^{(d)})$. In this case the assumption on f_g can be obtained from conditions on R_{g_0} , φ , F_g and $T_{st}^{g_0}$ in a similar fashion to what was done in the linear case (section 4), but we will not write them here for the sake of brevity.

Proof. Since K < 0, we can write K = -|K|. Similarly to the proof of proposition 5.1, we look for $v_- < v_+$ satisfying

$$\begin{split} L_g v_- + |K| v_-^{1-\frac{4}{d}} &\geq 0, \\ L_g v_+ + |K| v_+^{1-\frac{4}{d}} &\leq 0. \end{split}$$

Once again we take $v_{-} \equiv 0$. If v_{+} is constant and positive,

$$L_g v_+ + |K| v_+^{1-\frac{4}{d}} = f_g v_+ + |K| v_+^{1-\frac{4}{d}} \le -\min_M |f_g| v_+ + |K| v_+^{1-\frac{4}{d}},$$

where we used that $f_g < 0$ (notice that the minimum of $|f_g|$ will be non-zero by compactness of M). Setting the right hand side of above expression to be less than or equal to zero implies

$$v_+ \ge \left(\frac{|K|}{\min_M |f_g|}\right)^{\frac{a}{4}}$$

and hence we can set

$$v_+ \equiv \left(\frac{|K|}{\min_M |f_g|}\right)^{\frac{d}{4}}.$$

From proposition A.3 we obtain a non-negative solution v obeying the bound (5.4). This solution is strictly positive. In fact, since

(5.7)
$$L_g v = -|K| v^{1-\frac{4}{d}} \le 0$$

and $f_g < 0$, the maximum principle guarantees that v > 0 (unless v is constant).

To conclude this section, let us touch upon the boundedness of the effective potential for the special case of n = 2 and d > 4, stating and proving a version of conjecture 1.1. Indeed, by equation (2.41) in [14], we see that, the effective potential evaluated at a critical point is $\frac{(d-2)\alpha}{2dG_N}$. However, one cannot define a functional \mathcal{F} as we did earlier by simply plugging in a critical point. This is because, unlike the linear case d = 4, if there is more than one critical point, \mathcal{F} may not be well-defined.

This being said, in the case of n = 2, it can be proven that, regardless of which critical point is plugged into the effective potential, it is still bounded from below. Firstly, we see that the statement is non-trivial only when $\alpha < 0$. Hence, we shall assume that $\alpha = -|\alpha|$. Upon introducing a new function $u = v|\alpha|^{-\frac{d}{4}} > 0$, equation (5.1) becomes,

(5.8)
$$\Delta_g u - \frac{1}{2(d-1)} \left(\frac{d}{2}R_g + T^{(d)}\right) u = -\frac{d-2}{4(d-1)} u^{1-\frac{4}{d}}$$

Imposing the warped volume constraint, it follows that

$$\mathcal{V}(g) = \frac{d-2}{2dG_N} \frac{1}{\left(G_N \int_M u^{2-\frac{4}{d}} \, dV_g\right)^{\frac{2}{d-2}}}$$

Proposition 5.5. Let d > 4. Assume that there exists a constant $\eta > 0$ such that

$$\int_M T^{(d)} \, dV_{\widetilde{g}} \le \eta, \quad \text{for all} \quad \widetilde{g} \in [g].$$

Then there exists a number $K_{\eta} \in \mathbb{R}$ such that $\mathcal{V}(g = e^{2\varphi}g_0) \geq K_{\eta}$ for all smooth functions φ such that the volume of g is fixed.

Proof. As usual we will use C > 0 to denote several different constants independent of φ . Dividing equation (5.8) by u and integrating (with respect to dV_g), we have,

$$-\int_{M} \frac{d-2}{4(d-1)} u^{-\frac{4}{d}} dV_{g} = \int_{M} \frac{|\nabla u|^{2}}{u^{2}} dV_{g} - \frac{1}{2(d-1)} \left(\frac{d}{2} \int_{M} R_{g} dV_{g} + \int_{M} T^{(d)} dV_{g}\right) > -C.$$

where we have used that $\frac{d}{2} \int_M R_g dV_g$ is a topological invariant by the Gauss-Bonnet theorem, and hence independent of φ . Therefore,

(5.9)
$$\int_{M} u^{-\frac{4}{d}} dV_g < C$$

Let $w = u^{\frac{4}{d}}$. Then,

$$\mathcal{V}(g) = -\frac{C}{(\int_M w^{\frac{d}{2}-1} \, dV_g)^{\frac{2}{d-2}}} \ge -\frac{C}{\int_M w \, dV_g} \ge -C \int_M \frac{1}{w} \, dV_g.$$

The last inequality (which is obtained by Jensen's inequality) along with inequality (5.9) implies the result. $\hfill \Box$

6. Examples

Here, it is pointed out how our results can be applied in some important examples. As we do not intend to exhaust all possible models, we will keep the discussion short and somewhat informal, highlighting the main ideas and avoiding technicalities⁹.

As was stressed in section 2.1, important examples involving the string term are $T_{\rm st}^{g_0} = \Lambda \delta_{g_0}(p)$ and $T_{\rm st}^{g_0} = \Lambda \delta_{g_0}(N)$, for some $p \in M$, some function Λ which does not depend on the metric, and some submanifold N. In both cases, (2.14) holds, and an inspection in the proof of theorem 2.1-(2) reveals that this is essentially all that is needed for the proof to work. Indeed, from (2.14) it follows that (3.16) remains valid, and the rest of of the arguments go through. We conclude that theorem

⁹As we will be looking at specific examples, the metrics under consideration will be of the form $g = e^{2\varphi}g_0$ for some particular family of conformal factors φ . It is clear that we can still apply theorem 2.1 when the infimum is taken over a smaller set $\mathcal{P} \subset \mathcal{S}$.

(2.1) can also be applied when the string term is a delta function or an orientifold plane (provided, of course, that the remaining hypotheses are still in place).

In a similar fashion, suppose $T_{\rm st}^g$ takes a Gaussian shape,

(6.1)
$$T_{\rm st}^g = \frac{1}{\sigma_g^n} \exp\left(-\frac{r_g^2}{2\sigma_g^2}\right),$$

where r_g is the Riemannian distance to a fixed point $p \in M$, and the "width" σ_g is allowed to depend on g. For definiteness, we can imagine to be working in the neighborhood of p, so that (6.1) is multiplied by a suitably chosen cut-off function outside a ball centered at p.

In order to investigate how should σ_g transform, one can use the transformation law for a delta function,

(6.2)
$$\delta_g(p) = e^{-n\varphi} \delta_{g_0}(p), \text{ for } g = e^{2\varphi} g_0$$

as a guide. Upon rescaling of the metric,

$$g = \lambda^2 g_0$$

one obtains,

$$\int_{M} \frac{1}{\sigma_{g}^{n}} \exp\left(-\frac{\lambda^{2} r_{g_{0}}^{2}}{2\sigma_{g}^{2}}\right) \lambda^{n} dV_{g_{0}}.$$

If this is to equal

$$\int_{M} \frac{1}{\sigma_{g_{0}}^{n}} \exp\left(-\frac{r_{g_{0}}^{2}}{2\sigma_{g_{0}}^{2}}\right) \, dV_{g_{0}},$$

then we see that σ_g has to transform as

$$\sigma_g = \lambda \sigma_{g_0}.$$

Hence, requiring σ_q to transform as above reproduces the behavior (6.2) for rescalings of the metric.

For more general conformal transformations, $g = e^{2\varphi}g_0$, the distance function appearing in (6.1) does not change in a simple way, and hence finding the correct transformation law for σ_g is complicated. But when σ_{g_0} is very small, we hope for a behavior similar to (6.2), since in this case $T_{\rm st}^{g_0}$ will approach a delta function. We can then impose on σ_g a condition that guarantees a transformation law of the form

$$\frac{1}{\sigma_g^n} \exp\left(-\frac{r_g^2}{2\sigma_g^2}\right) = e^{-n\varphi} \frac{1}{\sigma_{g_0}^n} \exp\left(-\frac{r_{g_0}^2}{2\sigma_{g_0}^2}\right) + O(\sigma_0) \quad \text{as} \quad \sigma_0 \to 0.$$

In this situation, the integral of the string term will again be bounded independent of φ , and the arguments of theorem 2.1 may be used to prove that the effective potential is bounded from below.

Another example involving radially symmetric functions is given by

(6.3)
$$e^{2\varphi} = \frac{1}{(a^2 + |x|^2)^{\gamma}}$$

with a > 0 and $\gamma > 1$. This example has already been considered in [14], so let us investigate how it fits in the present results¹⁰. Again for definiteness, a suitable cut-off function has be to used far

¹⁰In [14], $\gamma < 2$ is also imposed. This is to ensure that the scalar curvature is not negative for large |x|, but we will not need this condition here.

away from the origin. More precisely, (6.3) is defined in \mathbb{R}^n , with |x| being the Euclidean distance to the origin, and we are considering the conformally flat metric

(6.4)
$$g_{\gamma} = g = e^{2\varphi}\delta = \frac{1}{(a^2 + |x|^2)^{\gamma}}\delta,$$

with δ being the Euclidean metric. Choosing an appropriate cut-off function for |x| large, this metric can be glued to the manifold (M, g_0) , as indicated in [14].

Remark 6.1. With the above gluing argument implicitly understood, we will work in \mathbb{R}^n for simplicity. In other words, we will explore the behavior of g in \mathbb{R}^n , but our primary interest is in the restriction of all quantities to a fixed ball of large radius.

Firstly, notice that the constraint $\int_M e^{n\varphi} dV_g = A$ fixes a in (6.3):

$$\int_{\mathbb{R}^n} \frac{dx}{(a^2 + |x|^2)^{\frac{n\gamma}{2}}} = a^{-n(\gamma-1)} \operatorname{vol}(S^{n-1}) \int_0^\infty \frac{s^{n-1}}{(1+s^2)^{\frac{n\gamma}{2}}} ds.$$

The condition $\gamma > 1$ ensures the convergence of the integral, so that

(6.5)
$$a = a(\gamma) = \left(\frac{\operatorname{vol}(S^{n-1})\int_0^\infty \frac{s^{n-1}}{(1+s^2)^{\frac{n\gamma}{2}}} ds}{A}\right)^{\frac{1}{n(\gamma-1)}}$$

Using (2.9) one finds

(6.6)
$$R_g(r) = (a^2 + r^2)^{\gamma} \left[2(n-1)r^{1-n} \left(\frac{\gamma n r^{n-1}}{a^2 + r^2} - \frac{2\gamma r^{n+1}}{(a^2 + r^2)^2} \right) - \frac{\gamma^2 (n-2)(n-1)r^2}{(a^2 + r^2)^2} \right]$$

where r = |x|. It then follows that

(6.7)
$$R(0) = 2(n-1)n\gamma a^{2(\gamma-1)}$$

From (6.5), we see that $a \to \infty$ when $\gamma \to 1$, and hence the scalar curvature blows up at the origin. On the other hand, a remains bounded when γ increases. Hence, if one fixes $\gamma_0 > 1$ and considers values of γ satisfying $\gamma_0 < \gamma$, the scalar curvature is bounded at the origin, and it is not difficult to see from (6.6) that it remains bounded on compact sets. Therefore, given $\gamma_0 > 1$ and $\rho > 0$, there exists a constant $K_{\gamma_0,\rho}$ such that

$$|R_{g_{\gamma}}|(x) \le K_{\gamma,\rho},$$

for all $\gamma > \gamma_0$ and $|x| \le \rho$, where g_{γ} is given by (6.4). Therefore theorem 2.1 can be applied (after suitable adjustments as pointed out in remark 6.1), and we conclude that the effective potential is bounded from below for this family of metrics.

The previous example clearly resembles a well known phenomena on the round sphere, which we now briefly recall (see [26] for details).

Identifying the complement of the north pole of S^n with \mathbb{R}^n via stereographic projection, the round metric can then be written as

$$g_0 = 4u_0^{\frac{4}{n-2}}\delta_s$$

where u_0 is the radially symmetric function

$$u_0(x) = (1 + |x|^2)^{\frac{2-n}{2}}$$

which is usually referred to as the "standard bubble". Acting on g_0 with the group of conformal diffeomorphisms of the sphere generates the family of metrics

$$g_{\varepsilon} = 4u_{\varepsilon}^{\frac{4}{n-2}}\delta,$$

where the one-parameter family of functions

$$u_{\varepsilon}(x) = \varepsilon^{\frac{n-2}{2}} (\varepsilon^2 + |x|^2)^{\frac{2-n}{2}}, \ \varepsilon > 0,$$

is also commonly known as the standard bubble. The presence of the parameter $\varepsilon > 0$ is easily understood: the functions u_{ε} , or equivalently the metrics g_{ε} , come from dilations in \mathbb{R}^n inducing conformal transformations on the sphere. All the metrics g_{ε} have the same volume, and all of them are of constant scalar curvature equal to the round one. Hence, the volume constraint and the bound on the scalar curvature required in theorem 2.1 are satisfied for this family of metrics, and therefore the theorem can be applied. The important point here is that even though one restricts to constant scalar curvature metrics on the sphere, there is an entire non-compact family of them, namely g_{ε} , a situation where one would naturally wonder if theorem 2.1 holds.

The previous example admits a generalization in the following sense. Suppose $n \geq 3$, fix a background metric g_0 , and consider the set of φ such that the metric $g = e^{2\varphi}g_0$ has fixed volume and constant scalar curvature. When the Yamabe invariant of M is non-positive there exists a unique such φ [37, 38, 39]. In the the positive case, however, it is possible to find a dense (in the C^0 topology) subset \mathcal{P} of such conformal factors satisfying the volume constraint and having constant scalar curvature [32]. As the scalar curvature remains unchanged for any $\varphi \in \mathcal{P}$, the hypotheses of theorem 2.1 are valid, and therefore the effective potential will be bounded from below when the infimum is taken over \mathcal{P} .

In the situation described in the previous paragraph, more can be said. Due to compactness theorems for the Yamabe problem ([23] and references therein), the $C^{k,\alpha}$ norm of $e^{2\varphi}$ will be uniformly bounded for all $\varphi \in \mathcal{P}$ by a constant depending only on the background metric g_0^{11} . Combining this with equation (2.15) and Schauder estimates (see proposition A.4) gives the following bound for the critical points:

$$|| u ||_{C^{k,\alpha}(M)} \le C_1 + C_2 || u ||_{C^0(M)},$$

where the constants C_1 and C_2 depend only the background metric g_0 and $C^{0,\alpha}$ -norms of F_{g_0} and $T_{\rm st}^{g_0}$ (and of course on the dimension and on α). Furthermore, if the solution u is unique, then the above estimate holds with $C_2 = 0$ (possibly after redefining C_1). The important point here is that the growth of the derivatives of u (and of u itself when the solution is unique) is essentially controlled by the *fixed* quantities F_{g_0} and $T_{\rm st}^{g_0}$.

APPENDIX A. CONVENTIONS AND SUMMARY OF FORMULAE

In this appendix we state some results we used in the proofs. Even though they are standard, they are included here for the reader's convenience. The definitions and theorems here presented are far from general, but they will suffice to the purpose of this paper. We will comment on such generalizations, referring to the literature for details. We will also indicate in this appendix our sign conventions.

¹¹Not to miss a subtle point, such theorems hold only when $n \leq 24$ and the manifold is not conformally equivalent to the round sphere, and they also rely on the validity of the Positive Mass theorem in higher dimensions. But none of such issues arises in the cases of interest: we have already discussed the case of the round sphere, and the Positive Mass theorem holds in dimensions less than or equal to seven [41, 42].

A.1. Inequalities. Here we recall some standard inequalities. They can all be found, for example, in [16].

Cauchy's inequality with ε .

$$ab \leq \frac{\varepsilon}{2}a^2 + \frac{1}{2\varepsilon}b^2, \ a, b > 0, \ \varepsilon > 0$$

Hölder's inequality. Assume $1 \le p, q \le \infty$, $\frac{1}{p} + \frac{1}{q} = 1$. Then if $u \in L^p(M)$, $v \in L^q(M)$, we have

$$\int_{M} |uv| \, dV \le \| \, u \, \|_{L^p(M)} \| \, v \, \|_{L^q(M)}$$

Interpolation inequality. Assume $1 \le s \le r \le t \le \infty$ and

$$\frac{1}{r} = \frac{\theta}{s} + \frac{1-\theta}{t}$$

Suppose that $u \in L^{s}(M) \cap L^{t}(M)$. Then $u \in L^{r}(M)$ and

$$|| u ||_{L^{r}(M)} \leq || u ||_{L^{s}(M)}^{\theta} || u ||_{L^{t}(M)}^{1-\theta}$$

Jensen's inequality. Assume $f : \mathbb{R} \to \mathbb{R}$ is convex, $u : M \to \mathbb{R}$ is integrable and M has finite measure. Then

$$f\left(\int_{M} u \, dV\right) \le \int_{M} f(u) \, dV$$

where $f_M = \frac{1}{\operatorname{vol}(M)} \int_M$.

A.2. **Theorems.** Throughout this section we consider elliptic linear partial differential operators which, in local coordinates, are written as^{12}

(A.1)
$$Pu = \partial_j (a^{ij} \partial_i u) + b^i \partial_i u - cu,$$

where the functions a^{ij} , b^i and c are assumed to be smooth, and manifolds are always assumed to be compact and without boundary.

A differential operator is elliptic when its symbol is invertible [22, 4]. For scalar equations, this corresponds to saying that the matrix a^{ij} in (A.1) is positive definite for all $x \in M$. In other words, there exists a constant $\lambda > 0$, called the *ellipticity constant*, such that

$$a^{ij}(x)\xi_i\xi_j \ge \lambda |\xi|^2$$

for any $x \in M$ and tangent vectors $\xi \neq 0$ (where the above expression should be understood in local coordinates around x). The standard model of an elliptic operator is the Laplacian of a metric g, Δ_g .

Recall that a linear system in finite dimensions, written as Ax = b, has a solution if, and only if, b is orthogonal to the kernel of A^* . Much of the power of elliptic theory stems from the fact that elliptic operators enjoy a similar property:

Proposition A.1. (Fredholm Alternative, [22, 16, 19]). Let M be a smooth Riemannian manifold and let $P: C^{\infty}(M) \to C^{\infty}(M)$ be a linear elliptic differential operator, and denote by P^* its formal adjoint. Given $f \in C^{\infty}(M)$, the equation

$$(A.2) Pu = f$$

 $^{^{12}\}mathrm{In}$ particular, it is also assumed that all operators are of second order.

has a solution if, and only if, f is L^2 orthogonal to the kernel of P^* . In particular, if $P^* = P$ then equation (A.2) always has a solution, which is unique and smooth, provided that P has trivial kernel.

Proposition A.1 can be generalized to include differential operators between spaces of sections of vector bundles over M, rough coefficients and situations of very low regularity. See the above references, [44] and references therein.

For self-adjoint operators, the Fredholm alternative reduces the problem of existence to that of uniqueness. Uniqueness can be tackled with the maximum principle:

Proposition A.2. (Maximum principle, [22, 16, 19]) Consider the operator P in (A.1) on a domain $\Omega \subseteq M$ (possibly $\Omega = M$).

(i) Assume c = 0 and $Pu \ge 0$ (resp. ≤ 0) in Ω . If u achieves its maximum (resp. minimum) in the interior of Ω then u is constant.

(ii) Assume $c \ge 0$ and $Pu \ge 0$ (resp. ≤ 0) in Ω . Then u cannot achieve a non-negative maximum (resp. non-positive minimum) in the interior of Ω unless it is constant.

(iii) Assume $c \ge 0$ with c not identically zero. If Pu = 0 in M, then $u \equiv 0$.

Again, this can be generalized to include rough coefficients. Such generalizations usually involve replacing the previous inequalities by some inequality to hold in the sense of distributions. See [19, 45] and references therein.

Proposition A.3. (Elliptic regularity). Let u be a solution of Pu = f.

(1) If $f \in C^{k,\alpha}(M)$, then $u \in C^{k+2,\alpha}(M)$. In particular u is smooth if f is smooth.

(2) If
$$f \in L_k^p(M)$$
, then $u \in L_{k+2}^p(M)$.

Elliptic regularity says that solutions always "gain two derivatives" as compared to the right hand side term. This leads to the following useful bootstrap argument. Consider a non-linear equation of the form

$$Pu = f(x, u),$$

where $f: M \times \mathbb{R} \to \mathbb{R}$ is smooth, and suppose that somehow one manages to produce a solution u which is in $C^{0,\alpha}$. Since the composition of a smooth function with a Hölder continuous function is again Hölder continuous, we have $f(\cdot, u) \in C^{0,\alpha}(M)$. Hence elliptic regularity tells us that $u \in C^{2,\alpha}(M)$. But then $f(\cdot, u)$ is in fact in $C^{2,\alpha}(M)$, and so u must be in $C^{4,\alpha}(M)$. Continuing this argument yields $u \in C^{\infty}(M)$. Of course, if f is, say, only C^k , we can iterate this argument only up to order k.

Proposition A.4. (Schauder estimates [22, 19, 4]) Let M be a smooth Riemannian manifold and let $P: C^{\infty}(M) \to C^{\infty}(M)$ be a linear elliptic differential operator. There exist a constant C, depending only on the α , the dimension of M, the $C^{0,\alpha}$ -norm of a^{ij} , b^i and c, and on the ellipticity constant λ , such that

$$\| u \|_{C^{k+2,\alpha}(M)} \le C(\| P u \|_{C^{k,\alpha}(M)} + \| u \|_{C^{0}(M)})$$

Moreover, if one restricts u so that it is orthogonal (in L^2) to the kernel of P, then we can drop the C^0 term on the right hand side (after replacing C with a new constant depending only on the same quantities as before).

The next proposition is a powerful tool to produce solutions to semi-linear equations, as it can be seen, for example, from its applications in section 5. Its weakness relies on the fact that it does not guarantee uniqueness of solutions (and in fact it is easy to construct examples where the proposition applies but solutions are not unique, see [22]).

Proposition A.5. (Sub- and super-solutions, [40, 22]) Let M be a smooth Riemannian manifold. Consider the semi-linear elliptic equations

(A.3) $\Delta_g u + f(x, u) = 0,$

where $f \in C^{\infty}(M \times \mathbb{R})$. Suppose there there exist $\phi, \psi \in C^2(M)$ satisfying

$$\Delta_g \phi + f(x, \phi) \ge 0,$$

$$\Delta_g \psi + f(x, \psi) \le 0,$$

(such ϕ and ψ are called respectively a sub-solution and super-solution for (A.3)), and $\phi \leq \psi$. Then (A.3) has a solution $u \in C^{\infty}(M)$ such that $\phi \leq u \leq \psi$.

Proposition A.5 admits further generalizations, including for manifolds with boundary; see [35, 12].

A.3. Notation and conventions. Some notation and conventions we use:

Laplacian: $\Delta_g = \frac{1}{\sqrt{|g|}} \partial_i (\sqrt{|g|} g^{ij} \partial_j)$. This differs by a sign from [14].

Spacetime metric: $- + + \cdots +$.

Dimension: spacetime before compactification: D, spacetime after compactification: d, compact dimensions: n.

 $L_k^p(M,g)$ denotes the Sobolev space of k^{th} weakly differentiable functions which belong, along with its weak derivatives, to $L^p(M,g)$, where the measure of integration is the volume element of the metric g. We sometimes write $L_k^p(M)$ for simplicity (in fact, when M is compact, Sobolev spaces defined using different metrics are equivalent, see for example [33]).

 $C^{k,\alpha}(M)$ denotes the Hölder spaces with k derivatives and Hölder exponent α .

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DISCONZI, DOUGLAS, AND PINGALI

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30