Critical Points and Supersymmetric Vacua, III: String/M Models*

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Abstract: A fundamental problem in contemporary string/M theory is to count the number of inequivalent vacua satisfying constraints in a string theory model. This article contains the first rigorous results on the number and distribution of supersymmetric vacua of type IIb string theories compactified on a Calabi-Yau 3-fold \(X\) with flux. In particular, complete proofs of the counting formulas in Ashok-Douglas [AD] and Deneef-Douglas [DD1] are given, together with van der Corput style remainder estimates.

Supersymmetric vacua are critical points of certain holomorphic sections (flux superpotentials) of a line bundle \(\mathcal{L} \to \mathcal{C}\) over the moduli space of complex structures on \(X \times T^2\) with respect to the Weil-Petersson connection. Flux superpotentials form a lattice of full rank in a \(2b_3(X)\)-dimensional real subspace \(S \subset H^0(\mathcal{C}, \mathcal{L})\). We show that the density of critical points in \(\mathcal{C}\) for this lattice of sections is well approximated by Gaussian measures of the kind studied in [DSZ1, DSZ2, AD, DD1].

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1. Introduction

This is the third in a series of articles [DSZ1, DSZ2] (see also [Ze2]) by the authors on statistics of critical points of random holomorphic sections and their applications to the vacuum selection problem in string/M theory. We recall that, in these articles, a 'vacuum' in string theory is a Calabi-Yau manifold of complex dimension $d = 3$ which forms the 6 'small dimensions' of the 10-dimensional universe, together with a choice of orientifolding and flux. Mathematically, vacua are critical points of a superpotential $W$, a holomorphic section of a line bundle $L \to \mathbb{C}$ over the configuration space $C$ which will be recalled in §1.1. The 'vacuum selection problem' is that there exists no principle at present which selects a unique superpotential, nor a unique critical point of a given superpotential, out of a large ensemble of possible vacua. This motivates the program of studying statistics of vacua, whose basic problems are to count the number of vacua satisfying physically natural constraints and to determine how they are distributed in $C$ (see [Do, DD1, AD, DGKT, KL, Si]). In this article, we present the first rigorous results on counting vacua with remainder estimates. In particular, we justify and improve on the approximations made in [DD1].

Our previous articles [DSZ1, DSZ2] were devoted to the statistics of critical points of Gaussian random holomorphic sections of line bundles over complex manifolds. The principal issue we face in this article is that the physically relevant ensembles of superpotentials are not Gaussian but rather are discrete ensembles of 'quantized flux' superpotentials which form a set of lattice points in a hyperbolic shell in $H^3(X, \mathbb{C})$. This hyperbolic shell is defined by the inequality (known as the tadpole constraint)

$$0 \leq Q[\varphi] \leq L.$$  \hspace{1cm} (1)
where

\[ Q[\varphi] = Q(\varphi, \bar{\varphi}) = -\sqrt{-1} \int_X \varphi \wedge \bar{\varphi} \]  

(2)

is the Hodge-Riemann bilinear form. As will be recalled in §2.4, \( Q \) is an indefinite quadratic form, whose ‘null cone’ \( \{ G : Q(G) = 0 \} \) is a real quadric hypersurface which separates \( H^3(X, \mathbb{C}) \)) into the interior \( \{ W : Q(G) > 0 \} \) and the exterior where \( Q(G) < 0 \). As will be seen below (Propositions 3.1 and 2.1), only flux superpotentials corresponding to lattice points in \( \{ G : Q(G) > 0 \} \) contribute vacua, and that is why we consider the shell (1).

Our main results show that as \( L \to \infty \), the statistics of critical points relative to the discrete lattice ensemble is well approximated by the statistics of critical points relative to the continuum ensemble in the shell, which is equivalent to the Gaussian ensembles of [DSZ1, DSZ2] and is therefore well understood. Thus, the vacuum statistics problem in string/M theory is a mixture of two kinds of equidistribution problems:

1. The distribution of radial projections of lattice points onto a quadric hypersurface;
2. The distribution of critical points of a continuous ensemble of random holomorphic sections (related to a Gaussian ensemble) of a negative line bundle, and their interpretation in the special geometry of Calabi-Yau moduli spaces.

The equidistribution problem in (2) is analyzed in detail in [DSZ1, DD1], so the main purpose of this paper is to analyze (1) and to combine it with the previous analysis of (2).

At the end of this article in §7, we compare the mathematical results of this article to discussions of vacua in the string theory literature.

1.1. Background to the results.

To state our results, we will need some notation (see §2 for more details). The models we consider in this article are called type IIb flux compactifications [GVW, GKP]. We fix a complex 3-dimensional Calabi-Yau manifold \( X \), i.e. a complex manifold with trivial canonical bundle \( K_X \simeq \mathcal{O} \) and with first Betti number \( b_1(X) = 0 \). In some of the physics literature, it is also assumed that \( H^2,0(X) = 0 \), but our results hold without this assumption. For each complex structure \( z \) on \( X \), there is a corresponding Hodge decomposition

\[ H^3(X, \mathbb{C}) = H^{3,0}_z(X) \oplus H^{2,1}_z(X) \oplus H^{1,2}_z(X) \oplus H^{0,3}_z(X). \]  

(3)

The space \( H^{3,0}_z(X) \) of (3, 0)-forms relative to \( z \) is one-dimensional and is spanned by a nowhere vanishing holomorphic volume form \( \Omega_z \). We also put \( b_3 = b_3(X) = \dim H^3(X, \mathbb{R}) \), \( h^{p,q} = h^{p,q}(X) = \dim_{\mathbb{C}} H^{p,q}(X) \). Thus, \( b_3 = 2(h^{2,1} + 1) \).

When we speak of vacua of string theory compactified on the Calabi-Yau space \( X \), we refer to classical vacua of the effective supergravity theory it determines. As discussed in [St2], the effective supergravity Lagrangian is derived by ‘integrating out’ or neglecting the massive modes (positive eigenvalues) of various operators. The data of effective supergravity consists of \( (\mathcal{C}, \mathcal{L}, W) \) where:

1. \( \mathcal{C} \) is the configuration space;
2. \( \mathcal{L} \to \mathcal{C} \) is a holomorphic line bundle.
3. the superpotential \( W \) is a holomorphic section of \( \mathcal{L} \).
In type IIB flux compactifications the configuration space is the moduli space of Calabi-Yau (Ricci flat Kähler) product metrics on $X \times T^2$. At this time of writing, the study of vacua in string theory is simplified by replacing the moduli space of Calabi-Yau metrics by the moduli space of complex structures on $X$ (see e.g. [Do, AD]). In the case where $h^{2,0}(X) = 0$, this is equivalent to fixing the Kähler class $[\omega] \in H^2(X, \mathbb{R})$ of the Calabi-Yau metrics. Hence we define the configuration space to be

$$C = \mathcal{M} \times \mathcal{E},$$

(4)

where $\mathcal{M}$ is the moduli space of complex structures on $X$ and where $\mathcal{E} = \mathcal{H}/SL(2, \mathbb{Z})$ is the moduli space of elliptic curves. Throughout this paper we identify $C = \mathcal{M} \times \mathcal{E}$ with a fundamental domain $D$ for the modular group $\Gamma$ in the Teichmüller space $Teich(X) \times \mathcal{H}$ of complex structures (see §2.1). For simplicity of exposition, we refer to restrictions to $D$ of holomorphic objects on $Teich(X) \times \mathcal{H}$ as holomorphic objects over $C$.

The line bundle $\mathcal{L}$ is defined to be the dual line bundle to the Hodge bundle $H^{3,0}(X) \otimes H^{1,0}(T^2) \to C$, where $T^2 = \mathbb{R}^2/\mathbb{Z}^2$. We give $C$ the Weil-Petersson Kähler form $\omega_{WP}$ induced from the Weil-Petersson metric on $\mathcal{L}$ (see §3.3). To be precise, $\mathcal{L}$ is a holomorphic line bundle over $Teich(X) \times \mathcal{H}$, and $W$ is a holomorphic section of $Teich(X) \times \mathcal{H}$. But as mentioned above, by holomorphic sections $W \in H^0(\mathcal{L}, C)$ we mean restrictions to $D$ of holomorphic sections of $H^0(Teich(X) \times \mathcal{H}, \mathcal{L})$.

Type IIB flux compactifications contain two non-zero harmonic 3-forms $F, H \in H^3(X, \mathbb{Z})$ which are known respectively as the RR (Ramond-Ramond) and NS (Neveu-Schwarz) 3-form field strengths. We combine them into a complex flux $G = F + i H \in H^3(X, \mathbb{Z} \oplus i \mathbb{Z})$. The parameter $\tau \in \mathcal{E}$ is known as the dilaton-axion and may be viewed as the period of $\omega_\tau = dx + \tau dy$ over the one-cycle dual to $dy$ in $T^2$. Given $G \in H^3(X, \mathbb{Z} \oplus \sqrt{-1}\mathbb{Z})$, physicists define the corresponding flux superpotential $W_G$ by:

$$W_G(z, \tau) = \int_X (F + \tau H) \wedge \Omega_\tau,$$

(5)

where $\Omega_\tau \in H^{3,0}(X)$. This is not well-defined as a function on $C$ since $\Omega_\tau$ and $\tau$ depend on a choice of frame. To be more precise, $G \in H^3(X, \mathbb{C})$ determines a section $W_G$ of the line bundle

$$\mathcal{L} = (H^{3,0}(X) \otimes H^{1,0}(T^2))^* \to Teich(X) \times \mathcal{H}$$

by making $G$ into the following linear functional on $H^3_{\tau}(X) \otimes H^{1,0}_{\tau}(T^2)$:

$$(W_G(z, \tau), \Omega_\tau \otimes \omega_\tau) = \int_{X \times T^2} (F \wedge dy - H \wedge dx) \wedge (\Omega_\tau \wedge \omega_\tau).$$

(6)

The map $G \to W_G$ defines an injective real (but not complex) linear map which embeds complex integral fluxes

$$H^3(X, \mathbb{Z} \oplus \sqrt{-1}\mathbb{Z}) \to H^0(\mathcal{M}, \mathcal{L})$$

(7)

as a lattice of rank $2b_3$ in $H^0(\mathcal{M} \times \mathcal{E}, \mathcal{L})$ which we call the lattice $S^\mathbb{Z}$ of integral flux superpotentials. The real span

$$S = \mathbb{R}S^\mathbb{Z} \subset H^0(\mathcal{M}, \mathcal{L})$$

(8)
of $S^Z$ is also important, and will be referred to as the space of flux superpotentials. We emphasize here that $S$ is not a complex vector space, nor are any of the associated spaces discussed below. We also use the (real-linear) map $G \mapsto W_G$ to regard $Q$ as a quadratic form on $S$, writing

$$Q[W_G] := Q[G] = -\sqrt{-1} \int_X G \wedge \overline{G} = 2 \int_X F \wedge H,$$

$$G = F + iH \in H^3(X, \mathbb{C}).$$  \hspace{1cm} (9)

The bundles $H^3_{z,0} \rightarrow M$ and $H^1_{\tau,0} \rightarrow E$ carry Weil-Petersson Hermitian metrics $h_{WP}$ defined by

$$h_{WP}(\Omega_1, \Omega_1) = e^{-K(z, \bar{z})} = i \int_X \Omega_z \wedge \overline{\Omega_z},$$  \hspace{1cm} (10)

and their associated Chern connections $\nabla_{WP}$. They induce dual metrics and connections on $L$. We denote the connection simply by $\nabla$.

### 1.2. Statement of the problem.

Given a flux superpotential $W$, there is an associated potential energy on $C$ defined by

$$V_W(Z) = |\nabla W(Z)|^2 - 3|W(Z)|^2.$$

(See [WB] for background on $V$.) By a vacuum we mean a critical point of $V(Z)$ on $C$. In this paper, we only study supersymmetric vacua, namely $Z \in C$ which are connection critical points in the sense that $\nabla W_P W(Z) = 0$. We denote the set of supersymmetric vacua of $W$ by

$$\text{Crit}(W) = \{Z \in C : \nabla W_P W(Z) = 0\}.$$  \hspace{1cm} (12)

Our goal is thus to count and determine the distribution of the supersymmetric vacua

$$\{\text{SUSY vacua}\} = \bigcup_{G \in S^Z : Q[G] \leq L} \text{Crit}(W_G)$$  \hspace{1cm} (13)

as $W_G$ varies over the lattice $S^Z$ within the hyperbolic shell (1). To determine the distribution of vacua, we introduce the incidence relation

$$\mathcal{I} = \{(W_G, Z) \in S \times C : \nabla W_G(Z) = 0\}.$$  \hspace{1cm} (14)

We shall view $C$ as a fundamental domain for the modular group $\Gamma$ in Teichmüller space (cf. §2). The incidence variety $\mathcal{I}$ is then a real $2m$-dimensional subvariety of $C \times S$ with the following diagram of projections:

$$\mathcal{I} \subset C \times S \xrightarrow{\pi} S$$  \hspace{1cm} (15)

The fiber $\pi^{-1}(W)$ is the set $\text{Crit}(W)$ of critical points of $W$ in $C$. Since $C$ is regarded as a fundamental domain in Teichmüller space, the map $\pi$ is not surjective: there exist $W$
with no critical points in $C$; hence $\pi(C)$ is a domain with boundary in $S$ (see §6.4.1). Critical points can move out of $C$ as $W$ varies in $S$. (There is a similar but more complicated theory of non-supersymmetric vacua [DD2].)

The fibers of $\rho$ are the subspaces

$$S_Z := \{W \in S : \nabla_W \rho(W) = 0\},$$

which play a crucial role in this article. They have the remarkable Hodge theoretic identifications,

$$S_{z, \tau} \equiv H^{2,1}_z(X) \oplus H^{0,3}_z(X) \quad \text{(Proposition 3.1).}$$

It then follows (see Proposition 3.2) that $\mathcal{I} \xrightarrow{\rho} C$ is a vector bundle (with fiber $\approx \mathbb{C}^{b_3/2}$) over a manifold with boundary. Another key point is that the restrictions of $Q$ to the fibers are always positive definite:

$$Q_{z, \tau} := Q|_{H^{2,1}_z(X) \oplus H^{0,3}_z(X)} \gg 0 \quad \text{(Proposition 2.1),}$$

i.e. $S_Z$ lies in the positive cone $\{Q(\phi, \phi) > 0\}$ of the indefinite quadratic (Hodge-Riemann) form (2) (cf. §2.4).

We now define the discriminant locus

$$\tilde{D} = \{(Z, W) \in \mathcal{I} : \det H^\mathcal{C} W(Z) = 0\}$$

of points $(Z, W) \in \mathcal{I}$ such that $Z$ is a degenerate critical point of $W$, where $H^\mathcal{C} W(Z)$ is the complex Hessian of $W$ at the critical point $Z$ as defined in (59)–(61). Equivalently, $\tilde{D}$ is the set of critical points of the second projection $\mathcal{I} \xrightarrow{\pi} S$ together with the singular points of $\mathcal{I}$. Its image $D = \pi(\tilde{D})$ under $\pi$ is the discriminant variety of superpotentials with degenerate critical points.

For each $W \in S \setminus \{0\}$, we define its distribution of (non-degenerate) critical points as the measure $C_W$ on $\mathcal{I} \setminus \tilde{D}$ defined by

$$\langle C_W, \psi \rangle = \sum_{Z \in \text{Crit}(W)} \psi(Z, W),$$

for $\psi \in C(\mathcal{I})$ such that $\rho(\text{Supp } \psi)$ is relatively compact in $C$ and $\text{Supp } \psi$ is disjoint from $\tilde{D}$. A more general definition of $C_W$ is

$$C_W = |\det H^\mathcal{C} W(Z)| \nabla W^* \delta_0$$

which will be discussed in §4.2. We make these assumptions on $\psi$ so that the sum on the right side is a finite and well-defined sum. Indeed, the pull back is not well-defined (without further work) on $\tilde{D}$. We will say more about $\tilde{D}$ after the statement of Theorem 1.4.

The basic sums we study are:

$$N_{\psi}(L) = \sum \{\langle C_N, \psi \rangle : N \in S^Z, Q[N] \leq L\}$$

$$= \sum \{\psi(Z, N) : (Z, N) \in \mathcal{I}, N \in S^Z, 0 \leq Q[N] \leq L\}.$$

For instance, when $\psi \equiv \chi_K$ is the characteristic function of a compact subset $K \subset \subset \mathcal{I} \setminus \tilde{D}$, $N_{\psi}(L)$ counts the total number of non-degenerate critical points lying over $\rho(K)$.
coming from all integral flux superpotentials with $Q[W] \leq L$. Physicists are naturally interested in counting the number of vacua with close to the observed values of the cosmological constant and other physical quantities, and hence would study sums relevant to such quantities. For instance, the cosmological constant of the theory defined by a vacuum $Z$ is the value $V(Z)$ of the potential there (see [DD1], §3.3). Thus, we may state the main problem of this paper:

**Problem 1.1.** Find the asymptotics and remainder for $\mathcal{N}_\psi(L)$ as $L \to \infty$.

As indicated above, this problem is very closely related to the pure lattice point problem of measuring the rate of uniform distribution of radial projections of lattice points onto the surface of a quadric hypersurface. More generally one could consider any smooth strictly convex set $Q \subset \mathbb{R}^n$ ($n \geq 2$) with $0 \in Q^\circ$. Associated to $Q$ is the norm $|X|_Q$ of $X \in \mathbb{R}^n$ defined by

$$Q = \{X \in \mathbb{R}^n : |X|_Q < 1\}.$$  

To measure the equidistribution of radial projections of lattice points to $\partial Q$, one considers the sums

$$S_f(t) = \sum_{k \in \mathbb{Z}^n \cap \mathbb{R}^n_+ \setminus \{0\}} f\left(\frac{k}{|k|_Q}\right), \quad \text{with } f \in C^\infty(\partial Q), \ t > 0. \quad (22)$$

The parallel lattice point problem is then

**Problem 1.2.** Find the asymptotics and remainder for $S_f(t)$ as $t \to \infty$.

### 1.3. Statement of the results.

In Theorem 5.1, we obtain a van der Corput type estimate for the lattice point Problem 1.2. For the critical point problem, we first give an elementary formula which is based on a trivial lattice counting estimate (which is useful since it is sometimes sharp), namely where the remainder term is simply a count of the cubes of the lattice which intersect the boundary. We denote by $\chi_{QZ}$ the characteristic function of the shell $\{W \in S_Z : 0 < Q[Z][W] < 1\}$.

**Proposition 1.3.** Suppose that $\psi = \chi_K$, where $K \subset \mathcal{I}$ such that $(Z, W) \in K \leftrightarrow (Z, rW) \in K$ for $r \in \mathbb{R}^+$. Assume further that $\rho(K)$ is relatively compact in $C$ and $\pi(\partial K)$ is piecewise smooth. Then

$$\mathcal{N}_\psi(L) = L^b_3 \left[ \int_{C \setminus S_Z} \psi(Z, W) \left| \det H^c W(Z) \right| \chi_{QZ}(W) dW \, d\text{Vol}_{WP}(Z) + O\left( L^{-1/2} \right) \right].$$

Here and in Theorem 1.4 below, $dW$ means the multiple of Lebesgue measure on $S_Z$ which gives the volume form for the positive-definite quadratic form $Q_Z = Q|_{S_Z}$. We note that the integral converges, since by (18), $\{Q[Z] \leq 1\}$ is an ellipsoid of finite volume.

It would be interesting to know if the remainder estimate is sharp for any domain $K \subset \mathcal{I}$. In the pure lattice point Problem 1.2, the corresponding ‘trivial estimate’ is sharp. For instance, consider the domain $K = S^{n-1}_+ \subset S^{n-1}$ formed by the northern hemisphere and put $\psi = \chi_K$. Then the remainder term

$$\sum_{k \in \mathbb{Z}^n, |k| \leq \sqrt{L}} \chi_K \left(\frac{k}{|k|}\right) - L^{\frac{n}{2}} \int_K f \, dA$$
reflects the concentration of projections of lattice points on the boundary $\partial S_{n-1}^{n-1}$, namely a great equatorial sphere. When the equator is defined by $x_n = 0$, the lattice points projecting over the equator are the lattice points in $\mathbb{Z}^{n-1} \subset \mathbb{R}^{n-1}$ and the number with $|k| \leq \sqrt{L}$ is of size $\sim L^{n-1}/2$. Analogously one may ask if there are domains $K \subset C$ along which critical points concentrate to the same maximal degree. Some evidence that the answer is ‘no’ will be presented in §4.1.

Our main result stated below is a much sharper van der Corput type asymptotic estimate of $N_\psi(L)$ as $L \to \infty$ for homogeneous test functions which vanish near the discriminant locus. Here, we say that a function $\psi \in C(I)$ is homogeneous of order $\alpha$ if

$$
\psi(Z, rW) = r^\alpha \psi(Z, W), \quad (Z, W) \in I, \ r \in \mathbb{R}^+.
$$

We consider homogeneous functions since they include (smoothed) characteristic functions as well as the cosmological constant (which is homogeneous of degree 2).

\textbf{Theorem 1.4.} Let $\psi \in C^\infty(I)$ be homogeneous of order $\alpha \geq 0$ and suppose that $\rho(\text{Supp} \, \psi)$ is a compact subset of $C$ and $\text{Supp} \, \psi \cap \tilde{D} = \emptyset$. Then

$$
N_\psi(L) = L^{b_3 + \alpha/2} \left[ \int_C \int_{S_Z} \psi(Z, W) \left| \det H^c W(Z) \right| \chi_{QZ}(W) \ dW \ d\text{Vol}_W p(Z) + O \left( L^{-\frac{2b_3}{2b_3+1}} \right) \right].
$$

It is reasonable to make the assumption $\text{Supp} \, \psi \cap \tilde{D} = \emptyset$, because degenerate critical points cannot be physically acceptable vacua in string/M theory. Indeed, the Hessian of $W$ at a critical point defines the ‘fermionic mass matrix’ of the theory, and a degenerate critical point would give rise to massless fermions which are not observed in physics. (See [WB] for definitions of the mass matrix.)

Let us note some key features of the geometry of $\tilde{D}$ which play a role in the assumptions (and proofs) of Proposition 1.3 and Theorem 1.4. First, as observed in [DSZ1, DSZ2], its defining equation

$$
\det H^c W(Z) = \det(H^*H - |W|^2 I) = 0
$$

is real valued; here, $H$ is the holomorphic Hessian (see §3.2). Hence, $\tilde{D} \subset C$ is a real analytic hypersurface (with boundary). For test functions $\psi$ which do not vanish on $\tilde{D}$, the expression $\langle C_W, \psi \rangle$ (when well-defined) can jump as one passes from one component of $S \setminus D$ to another or across the boundary of $C$. It follows from (23) that $\tilde{D} \cap (\{Z\} \times S_Z)$ is a real conic hypersurface for all $Z \in C$. Thus $\tilde{D} \to C$ is a bundle of conic hypersurfaces and $\rho(\tilde{D}) = C$; i.e., every point of moduli space is a degenerate critical point of some superpotential. We further note that $S \setminus D$ consists of a finite number of connected components, and that $\pi : I \setminus \tilde{D} \to \pi(S) \setminus D$ is a finite covering over each connected component of $\pi(S) \setminus D$.

\textbf{1.4. Special geometry and critical point density.} In obtaining reliable order of magnitude results on numbers of vacua in a given string/M model, it is important to estimate the size of the leading coefficient

$$
\int_C \psi(Z) \int_{S_Z} \left| \det H^c W(Z) \right| \chi_{QZ}(W) \ dW \ d\text{Vol}_W p(Z)
$$
and of the remainder. Since little is known about the volume of $\mathcal{C}$ at present (cf. [LS1]), we concentrate on estimating the integrand

$$K_{\text{crit}}(Z) := \int_{S_Z} |\det H^c W(Z)| \chi_{Q_Z} dW$$

in the $b_3$ aspect. It is also important to study the behavior of the $K_{\text{crit}}(Z)$ as $Z$ tends to ‘infinity’ in $\mathcal{C}$, or to a singular point such as a conifold point (when one exists).

A key feature of $K_{\text{crit}}(Z)$ is that it is the integral of a homogeneous function of order $b_3$ over a space of dimension $\dim_{\mathbb{R}} S_Z = b_3 = 2(h^{2,1} + 1)$. Among the known Calabi-Yau 3-folds it is common to have $300 < b_3 < 1000$, hence the integral is often over a space of large dimension. The $b_3$-dependence is sensitive since (e.g.) the ratio of the $L^\infty$ norm to the $L^2$ norm of a homogeneous function of degree $b_3$ in $b_3$ variables can be of order $b_3^{-1}$. It is useful to have alternative formulas for the leading coefficient, and we now present a few. We will use them to suggest conjectures on the order of magnitude of $K_{\text{crit}}(Z)$ in the $b_3$ aspect in §7.

First, using the homogeneity of the integrand, we may rewrite the integral in terms of a Gaussian density

$$K_{\text{crit}}(Z) = \frac{1}{b_3!} \int_{S_Z} |\det H^c W(Z)| e^{-\langle Q_Z W, W \rangle} dW.$$  

This formula shows that $K_{\text{crit}}$ is formally analogous to density of critical points of random holomorphic sections relative to a Gaussian measure studied in [DSZ1]. For this reason, we call (24) the critical point density. However, the measure $e^{-Q[W]} \chi_{[0<Q<1]}(W)dW$ is of infinite volume, so the analogy should not be taken too literally. The density $K_{\text{crit}}(Z)$ is well-defined despite the infinite volume of the underlying measure on $S$ because the fibers $Q_Z$ of $\rho|_Q$ are of finite volume. Indeed, the conditional measures of $e^{-Q[W]}dW$ are standard (un-normalized) Gaussian measures $e^{-Q_Z(W)}dW$.

Next, we rewrite the integrals by the methods in [DSZ1, DSZ2]. The first method is to change variables to the Hessian $H^c W(Z)$, i.e. to ‘push-forward’ the $S_Z$ integral under the Hessian map

$$H_Z : S_Z \to \text{Sym}(m, \mathbb{C}) \oplus \mathbb{C}, \quad H_Z(W) = H^c W(Z),$$

where $m = \dim \mathcal{C} = h^{2,1} + 1$. In [DSZ1, DSZ2], we used this change of variables to simplify the formulas for the density of critical points. There, however, the spaces of holomorphic sections of the line bundles $L \to M$ were so large that the image of the Hessian map was the entire space $\text{Sym}(m, \mathbb{C}) \oplus \mathbb{C}$ of complex Hessians of rank equal to the dimension $m = \dim M$. In the case of type IIb flux compactifications, the dimension of the configuration space $\mathcal{C}$ is as large as the dimension of the space $S$ of sections, and the Hessian map is by no means surjective. Indeed, in Lemma 6.1, we prove that the Hessian map is an isomorphism to a real $b_3$-dimensional space $\mathcal{H}_Z \oplus \mathbb{C}$, where $\mathcal{H}_Z$ is spanned (over $\mathbb{R}$) by the $2h^{2,1}$ Hermitian matrices

$$\xi^j := \begin{pmatrix} e_j \\ e_j^* \end{pmatrix}, \quad \xi^{h^{2,1}+j} := \begin{pmatrix} e_j & ie_j \\ ie_j^* & -i \mathcal{F}^j(z) \end{pmatrix}, \quad j = 1, \ldots, h^{2,1}. \tag{27}$$

Here, $e_j$ is the $j^{th}$ standard basis element of $\mathbb{C}^{h^{2,1}}$ and $\mathcal{F}^j(z) \in \text{Sym}(h^{2,1}, \mathbb{C})$ is the matrix $(\mathcal{F}^j_{ik}(z))$ whose entries define the ‘Yukawa couplings’ on $\mathcal{M}$ (see (46), §2.3 or [St1, CO]) with respect to normal coordinates at the point $z \in \mathcal{M}$. 

Since $\mathcal{H}_Z$ is not a complex subspace of $\text{Sym}(m, \mathbb{C})$, we regard $\text{Sym}(m, \mathbb{C})$ as a real vector space with inner product

$$(A, B)_\mathbb{R} = \text{Re} \langle A, B \rangle_{HS} = \text{Re} \langle \text{Trace } AB^* \rangle. \quad (28)$$

To state our next result, we let $\Lambda_Z$ be the operator given by the distortion under the Hessian map (see §6.2):

$$((\Lambda_Z \oplus I_\mathbb{C})^{-1} H_Z W, H_Z W)_\mathbb{R} = Q[W] \quad (W \in S_Z),$$

where $Q[W]$ is given by (9). In terms of the basis $\{\xi^a\}_{1 \leq a \leq 2h^2, 1}$,

$$\Lambda_Z \xi^a = \sum_{b=1}^{2h^2, 1} \Lambda_{ab} \xi^b, \quad \Lambda_{ab} = (\xi^a, \xi^b)_\mathbb{R}. \quad (29)$$

The $\Lambda$ matrix has the block form

$$(\Lambda_{ab}) = \begin{pmatrix} \Lambda' & \Lambda'' \\ \Lambda' & \Lambda' \end{pmatrix}, \quad \Lambda'_{jk} = 2\delta_{jk} + \text{Re } \text{Tr } F^j F^{k*},$$

$$\Lambda''_{jk} = \text{Im } \text{Tr } F^j F^{k*}. \quad (30)$$

In Proposition 6.2, we show that the $(1, 1)$ form

$$\omega_{\Lambda} : = \frac{i}{2} \sum (\Lambda'_{jk} + i \Lambda''_{jk}) dz^j \wedge d\bar{z}^k$$

$$= \frac{i}{2} \sum [2\delta_{jk} + \text{Tr } F^j (z_0) F^{k*}(z_0)] dz^j \wedge d\bar{z}^k \quad (31)$$

is the so-called Hodge metric $(m + 3)\omega_{WP} + \text{Ric}(\omega_{WP})$ of the Weil-Petersson metric [Lu, Wa2].

By the injectivity of the Hessian map (stated in Lemma 6.1), we can make the change of variables $W^H_z(H, x)$ in (24)–(25) to obtain the following alternate formulas for $K_{\text{crit}}(Z)$:

$$K_{\text{crit}}(Z) = \frac{1}{\sqrt{\det \Lambda_Z}} \int_{\mathcal{H}_Z \oplus \mathbb{C}} \left| \det H^* H - |x|^2 I \right| \chi_{\Lambda_Z}(H, x) dH dx,$$

$$= \frac{1}{b_3! \sqrt{\det \Lambda_Z}} \int_{\mathcal{H}_Z \oplus \mathbb{C}} \left| \det H^* H - |x|^2 I \right| e^{-(\Lambda_Z^{-1} H, H)_{\mathbb{R}} - |x|^2} dH dx,$$

$$\quad (32)$$

where $\chi_{\Lambda_Z}$ is the characteristic function of the ellipsoid $\{\Lambda_Z^{-1} H, H\}_{\mathbb{R}} + |x|^2 \leq 1$. These formulas are analogous to Theorem 1 and Corollary 2 of [DSZ1], the key difference being that here we integrate over a moving subspace $\mathcal{H}_Z$ of symmetric matrices.

We similarly have the following alternative formulations of Proposition 1.3 and Theorem 1.4:
Corollary 1.5. Let $\psi = \chi_K$, where $K \subset T$ is as in Proposition 1.3, and let $\tilde{\psi}(Z, H_Z W) = \psi(Z, W).$ Then,

$$N_\psi(L) = \frac{L^{b_3}}{b_3!} \left[ \int_{T} \frac{1}{\sqrt{\det \Lambda_Z}} \int_{H_Z \oplus T} \tilde{\psi}(Z; H, x) \det H^* H \right.$$

$$- |x|^2 I \left. \right| e^{-(\Lambda_Z^{-1} H, H)_{R} - |x|^2} dH dx d\text{Vol}_{WP}(Z)$$

$$+ O(L^{-1/2}) \right].$$

Corollary 1.6. Let $\psi \in C^\infty(T)$ be homogeneous of order $\alpha \geq 0$ and suppose that $\rho(\text{Supp } \psi)$ is a compact subset of $T$ and $\text{Supp } \psi \cap \bar{D} = \emptyset.$ Let $\tilde{\psi}(Z, H_Z W) = \psi(Z, W).$ Then,

$$N_\psi(L) = \frac{L^{b_3+\alpha/2}}{\Gamma(b_3 + \alpha/2 + 1)} \left[ \int_{T} \frac{1}{\sqrt{\det \Lambda_Z}} \int_{H_Z \oplus T} \tilde{\psi}(Z; H, x) \right.$$

$$\times \left| \det H^* H - |x|^2 I \right| e^{-(\Lambda_Z^{-1} H, H)_{R} - |x|^2} dH dx d\text{Vol}_{WP}(Z)$$

$$+ O \left( L^{-\frac{2b_3}{2b_3+1}} \right) \right].$$

It is not obvious how to estimate the dependence of the integral for $K_{\text{crit}}(Z)$ on the subspace $H_Z$. There are two natural ways to parameterize this space. One (which is used in [DD1]) is to use as a basis of $H_Z$ the Hessians of a $Q_Z$-orthonormal basis of $S_Z$. A second method is to use the orthonormal basis of eigenmatrices $\{H_j\}$ of $\Lambda_Z$ with respect to the inner product (28). We thus put $\Lambda_Z H_j(Z) = \mu_j(Z) H_j(Z),$ and $H(y, Z) = \sum_j y_j H_j(Z).$ We also let $D(\mu)$ denote the diagonal matrix with entries $\mu_j.$ Changing variables to $\mu_j^{1/2} y$ cancels $\frac{1}{\sqrt{\det \Lambda_Z}}$ and we obtain:

Corollary 1.7. We have:

$$K_{\text{crit}}(Z) = \int_{|y|^2 + |x|^2 \leq 1} \left| \det H(D(\mu) y, Z)^* H(D(\mu) y, Z) - |x|^2 I \right| dy dx.$$

In §7 we will discuss some conjectural bounds on the density of critical point based on the assumption that the subspaces $H_Z$ are sufficiently random subspaces of $\text{Sym}(h^{2,1}, \mathbb{C}).$

1.5. Index density. The absolute value in the expressions for the distribution of critical points $C_W$ of a single section (20) and the expected distribution of critical points of a random section (e.g., (32)) make it very difficult to estimate the order of magnitude in $b_3$ of the density of critical points. A simplifying ‘approximation’ is to drop the absolute value around the determinant. The resulting density is the index density for critical points. It was used in [AD] and [DD1] to give a lower bound for the critical point density.

To be precise, we modify (20) by defining the signed distribution of critical points of $W$ as the measure $C_W$ on $\mathcal{I} \setminus \bar{D}$ given by

$$\langle \text{Ind}_W, \psi \rangle = \sum_{Z \in \text{Crit}(W)} \left( \text{sign } \det D^2 W(Z) \right) \psi(Z, W),$$

(33)
where sign $a = 1, 0, -1$ if $a$ is positive, 0, or negative, respectively. We then study the sums

$$\text{Ind}_\psi(L) = \sum \{ (\text{Ind}_N, \psi) : N \in S_Z, Q[N] \leq L \}. \quad (34)$$

For instance, if $\psi(Z, W) = \chi_K(Z)$ is the characteristic function of a compact set $K \subset C$, then $\text{Ind}_\psi(L)$ is the sum $\sum_{Z \in \text{Crit}(W) \cap K} (\text{sign det } D^2 W(Z))$ over all non-degenerate critical points lying over $K$ of all integral flux superpotentials with $Q[W] \leq L$.

Simultaneously with Proposition 1.3, we obtain formula (1.5) of Ashok-Douglas [AD] with an estimate for the error produced by passing from the sum to the integral (cf. §4):

**Theorem 1.8.** Let $K$ be a compact subset of $C$ with piecewise smooth boundary. Then

$$\text{Ind}_{\chi_K}(L) = \frac{(\pi L)^{b_3}}{b_3! 2^{b_3/2}} \left[ \int_K c_m(T^{*(1,0)}(C) \otimes L, \omega_{WP} \otimes h_{WP}^*) + O \left( L^{-1/2} \right) \right],$$

where $m = \dim C = b_3/2$ and $c_m(T^{*(1,0)}(C) \otimes L, \omega_{WP} \otimes h_{WP}^*) = \frac{1}{\pi} \det(-R - \omega_{WP} \otimes I)$ is the $m^{th}$ Chern form of $T^{*(1,0)}(C) \otimes L$ with respect to the Weil-Petersson metric $\omega_{WP} \otimes I$.

Here, $R = \sum_{i,j} R_{i,j}^k dz^i \wedge d\bar{z}^j$ is the curvature $(1, 1)$ form of $T^{*(1,0)}(C)$ regarded as an $m \times m$ Hermitian-matrix-valued 2-form (with $m = \dim C = b_3/2$) and $\omega_{WP} \otimes I$ is a scalar 2-form times the $m \times m$ identity matrix. The determinant is defined as in Chern-Weil theory. The only additional step in the proof is the evaluation (given in Lemma 6.3) of the analogue of (25) in terms of the curvature form:

$$\int_{S_Z} \det H^c W(Z) e^{-\langle Q Z W, W \rangle} dW = \left( \frac{\pi}{2} \right)^m \frac{\det(-R - \omega_{WP} \otimes I)}{d\text{Vol}_{WP}}. \quad (35)$$

Recall that the Chern-Gauss-Bonnet theorem tells us that if $W$ is a holomorphic section of a complex line bundle $L \to M_m$ over a compact complex manifold such that $\nabla W$ has only non-degenerate zeros, then

$$c_m(T^{*(1,0)} M \otimes L) = \text{Ind} \nabla W := \sum_{p : \nabla W(p) = 0} \text{sign } \det H^c W(p).$$

However, the Chern-Gauss-Bonnet theorem does not apply in our setting, and indeed $\text{Ind} \nabla W$ is not constant in $W$, since $C$ is an incomplete Kähler manifold and critical points can occur on the boundary or disappear. There exists a Chern-Gauss-Bonnet theorem for manifolds with boundary which expresses $\text{Ind} \nabla W$ as $c_m(E)$ plus a boundary correction depending on $W$, but the correction term involves integrating a differential form over the boundary and that becomes problematic when the boundary is highly irregular as in the case of $C$. Nevertheless, the theorem shows that asymptotically the average index density equals the Chern-Gauss-Bonnet form.
1.6. Relations to prior results in the physics and mathematics literature. We now relate our results to the physics literature on the number of vacua and the complexity of the string theory landscape as well as to the mathematical literature on lattice points. A more detailed discussion of the landscape aspects is given in §7.

First, the string/M aspects. Over the last five years or so, many physics articles have been devoted to estimating the number of candidate vacua $N_{vac}$ of string/M theory, in particular those which are consistent with the standard model. The candidate vacua are often pictured as valleys in a ‘string theory landscape’, which is the graph of the effective potential. The number of vacua is often stated as being around $10^{500}$. In [BP] Bousso-Polchinski related the number of vacua to the number of quantized fluxes $N$ satisfying a constraint $|N| \leq L$, which implies $N_{vac}(L) \sim \frac{b_3}{b_3^3}$ (see also [AD, Si]). In the specific type IIb flux compactifications studied in this paper, the constraint is hyperbolic rather than elliptic (as imagined in [BP]), and the more precise estimate $N_{vac}(L) \sim \frac{L b_3}{b_3^3} f(b_3)$ was given in [AD, DD1], where $f(b_3)$ is the moduli space integral of the Gaussian integral in (32); it will be discussed further in §7. There we will also review the heuristics and the mathematics of the landscape in more detail.

What do our results imply about the number of vacua? Since Proposition 1.3 and Theorem 1.4 are asymptotic results as $L \to \infty$, they are most useful when $L b_3$ is very large. But it is difficult to quantify ‘very large’ due to the complexity of the leading coefficient (24), of the remainder and of the volume of $C$. Hence, we cannot make precise estimates on the number of vacua at this time.

However, to bridge our results with estimates in string theory, we make a speculative attempt in §7.3 to draw order of magnitude conclusions from Theorem 1.4. We will use the symbol $\simeq$ in an informal sense of ‘same order of magnitude’ in $b_3$ (factorial, exponential and so on). There we give a heuristic estimate of $K_{\text{crit}}(Z) \simeq \frac{1}{b_3^2} (b_3/2)! \mu b_3^3$ for certain $\mu > 0$. More precisely, we give heuristic upper and lower bounds with different $\mu$ which are irrelevant when comparing factorials. To obtain an order of magnitude for $\frac{f(b_3)}{b_3^3}$ one would need to integrate $K_{\text{crit}}$ over $C$. At this time, the order of magnitude of the Weil-Petersson volume $Vol_{WP}(C)$ of $C$ is not known, even approximately (Z. Lu). We can however make a plausible estimate for the integral of $K_{\text{crit}}$ over the region where the norm of $Λ_Z$ is bounded by a uniform constant (independent of $b_3$). Since $Λ_Z$ is essentially the Hodge metric, regions where $|Λ_Z| \leq \mu$ are regions $K_\mu$ where the norm of the Ricci curvature of $ω_{WP}$ is bounded above by a uniform constant. It appears likely that the volume of such regions is bounded above by the volume of balls in $C^{b_3/2}$ of fixed radius (Z. Lu). Since the volume of balls in $C^{b_3/2}$ decays like $\frac{1}{(b_3/2)!}$, the number of vacua in $K_\mu$ would be approximately $\frac{L b_3}{b_3^3} \mu b_3^3$.

Now, in the physical models, $L$ is not a free parameter but is determined by $X$. In the case when there exists an involution $g$ of $X$ (an ‘orientifolding’) and a Calabi-Yau 4-fold $Z$ which is an elliptic fibration over $X/g$, the ‘tadpole’ number is then given by:

\[ \text{tadpole number} : \quad L = \chi(Z)/24. \]  

(36)

In many known examples [KLRY], one has $300 < b_3 < 1,000$ and $L \simeq C b_3$, where $1/3 \leq C \leq 3$. Hence the number of vacua in $K_\mu$ (and possibly in all of $C$) with the tadpole constraint $L \sim C b_3$ would have exponential growth $\frac{(C b_3)^{b_3}}{b_3!} \mu b_3^3$.

Next we turn to the purely lattice point aspects of the problem. From a mathematical point of view, this article combines statistical algebraic geometry in the sense of [BSZ1, DSZ1, DSZ2] with the study of radial projections of lattice points. As far as we know,
the radial projection of lattice points problem has not been studied systematically before in mathematics (we thank B. Randol for helping to sort out the historical background on this problem). The much harder problem of equidistribution of lattice points of fixed height $R$, i.e. lying on a sphere or hyperboloid of fixed radius $R$, has been studied by Yu. Linnik, C. Pommerenke [Po], W. Duke and others. But the remainders obtained in this more delicate problem are not as accurate as ours are for the bulk problem of projecting all lattice points of height $< R$. Counting projections of lattice points in domains of a hypersurface is equivalent to counting lattice points in certain cones, and there are some additional studies of this by methods of automorphic forms. In certain right circular cones with a flat top, Duke and Imamoglu [DO] use Dirichlet series and Shimura lifts to obtain the leading order asymptotics. Radial projections of lattice points additionally bear some resemblance to rational points. Some results and references for this problem are contained in [DO]. In [Ze1], the general problem of counting radial projections of lattice points in smooth domains of non-degenerate hypersurfaces is studied. In [NR], some further results are given on radial projections of lattice points, in particular in the case of hypersurfaces with flat spots or in the case of polyhedra.

2. Background on Calabi-Yau Manifolds and Complex Geometry

As mentioned in the introduction, the supersymmetric vacua of type IIB flux compactifications on a $CY_3$ are critical points of holomorphic sections of the holomorphic line bundle $L \to \mathcal{C}$ dual to the Hodge bundle $H^{4,0}(X \times T^2)$, where the configuration space $\mathcal{C}$ is the moduli space $\mathcal{M} \times \mathcal{E}$ of product complex structures on $X \times T^2$. In this section, we give the geometric background necessary for the analysis of critical points and Hessians of the holomorphic sections $W_G$ of (5).

The most significant aspects of Calabi-Yau geometry in the study of critical points of flux superpotentials are the following:

- The space $S_Z$ of flux superpotentials with $\nabla W_G(Z) = 0$ may be identified with the space $H^2_{\tau}(X)$ of fluxes $G = F + iH$ with the special Hodge decomposition $F + \tau H \in H^{2,1}_{\tau}(X) \oplus H^{0,3}_{\tau}(X)$. See Proposition 3.1.
- The space $H^2_{\tau}(X) \oplus H^{0,3}_{\tau}(X) \subset H^3(X, \mathbb{C})$ is a positive complex Lagrangian subspace. See Proposition 2.1. Hence, $S_Z$ is endowed with an inner product.

In addition, we review the relation between holomorphic derivatives, covariant derivatives and second fundamental forms for holomorphic frames $\Omega_1$ of the Hodge bundle, and recall the definition of the prepotential. These results are needed for the calculations in Lemmas 107 and 6.1. Much of this material is essentially standard [CO, St1, DD1], but it is not always stated precisely in the physics sources. We therefore present it here for the sake of clarity and completeness.

2.1. Geometry of Calabi-Yau manifolds. We recall that a Calabi-Yau $d$-fold $M$ is a complex $d$-dimensional manifold with trivial canonical bundle $K_M$, i.e. $c_1(M) = 0$. By the well-known theorem of Yau, there exists a unique Ricci flat Kähler metric in each Kähler class on $M$. In this article, we fix the Kähler class, and then the Calabi-Yau metrics correspond to the complex structures on $M$ modulo diffeomorphisms. We denote the moduli space of complex structures on $M$ by $\mathcal{M}$.

As mentioned in the introduction, the Calabi-Yau manifolds of concern in this article are the 4-folds $M = X \times T^2$, where $T^2 = \mathbb{R}^2/\mathbb{Z}^2$. The $T^2$ factor plays a special role, and
the geometric aspects mainly concern $X$. We only consider product Calabi-Yau metrics and complex structures on $M$. Thus, the configuration space $\mathcal{C} = M \times \mathcal{E}$, where $M$ is the moduli space of complex structures on $X$ and where $\mathcal{E}$ is the moduli space of elliptic curves. We denote a point of $\mathcal{C}$ by $Z = (z, \tau)$, where $z$ denotes a complex structure on $X$ and where $\tau$ denotes the complex structure on $T^2$ corresponding to the elliptic curve $\mathbb{C}/\mathbb{Z} \oplus \mathbb{Z} \tau$.

It is often simplest to view the moduli space of complex structures on $X$ as the quotient by the mapping class group $\Gamma$ of the Teichmüller space $Teich(X)$, where

$$Teich(X) = \{\text{complex structures on } X\}/\text{Diff}_0,$$

where $J \sim J'$ if there exists a diffeomorphism $\varphi \in \text{Diff}_0$ isotopic to the identity satisfying $\varphi^* J' = J$. The mapping class group is the group of connected components of the diffeomorphism group,

$$\Gamma_X := \text{Diff}(X)/\text{Diff}_0(X).$$

We shall identify $M$ with a fundamental domain for $\Gamma_X$ in $Teich(X)$, and $\mathcal{E}$ with the usual modular domain in $\mathcal{H}$.

The mapping class group for a Calabi-Yau $d$-fold has a representation on $H^d(M, \mathbb{R})$ which preserves the intersection form $Q$, which is symplectic in odd dimensions, and indefinite symmetric in even dimensions. In odd dimensions, this representation gives a homomorphism $\varphi : \Gamma_M \to \text{Sp}(b_d(M), \mathbb{Z})$, while in even dimensions it gives a homomorphism to the corresponding orthogonal group. It was proved by D. Sullivan [Sul] that if $d \geq 3$, then $\varphi(\Gamma_M)$ is an (arithmetic) subgroup of finite index (in $\text{Sp}(b_d(M), \mathbb{Z})$ if $d$ is odd), and the kernel of $\varphi$ is a finite subgroup.

On any CY manifold $M$ of dimension $d$, the space $H^{d,0}_z(M)$ of holomorphic $(d, 0)$ forms for a complex structure $Z$ is one-dimensional. It depends holomorphically on $Z$ and hence defines a complex holomorphic line bundle $L^*_M = H^{d,0} \to M$, which we refer to as the ‘Hodge bundle.’ The Hodge bundle is equipped with the Weil-Petersson (WP) Hermitian metric of (10), which we repeat here:

$$h_{WP}(\Omega, \Omega) = \frac{i}{2} \int_X \Omega \wedge \overline{\Omega}. \quad (37)$$

For a holomorphic Hermitian line bundle $(L, h) \to M$ and local holomorphic frame $e_L$ over an open set $U \subset M$, we write

$$|e_L(z)|^2_h = e^{-K(z)}. \quad (38)$$

The connection 1-form in this frame is the $(1, 0)$ form $-\partial K(z)$, and the curvature $(1, 1)$-form is given by

$$\omega = \frac{i}{2} \Theta_h = \frac{i}{2} \delta \partial K, \quad K = - \log |e_L|^2_h.$$

The Hermitian line bundle is said to be positive if $\omega$ is a positive $(1, 1)$ form, in which case $K$ is called the Kähler potential. The Hermitian line bundle $(L, h)$ is negative if $\omega$ is a negative $(1, 1)$ form.
In particular, the curvature of the Weil-Petersson metric on \( H^{d,0} \to \mathcal{M} \) is a positive \((1, 1)\) form on \( \mathcal{M} \), and hence it defines a Kähler form with potential (with respect to the frame \( \Omega_z \))

\[
K_{WP} = -\log h_{WP}(\Omega_z, \Omega_z) = -\log i \int_X \Omega \wedge \overline{\Omega}.
\] (39)

For instance, consider the Hodge bundle \( H^1_{\tau,0} \to E \). It has a standard frame \( d\tau + \tau d\bar{\tau} \) for which \( K = -\log \Im \tau \). Here, \( \tau \) is the standard coordinate on the upper half plane. Then \( \partial K = -\frac{i}{\tau - \bar{\tau}} d\tau \) and the Kähler form is \( -\frac{1}{2(\tau - \bar{\tau})^2} d\tau \wedge d\bar{\tau} \gg 0 \).

In the product situation of \( M = X \times T^2 \), \( H^{4,0}_{\tau,0}(X \times T^2) = H^{3,0}_\tau(X) \otimes H^{1,0}_0(T^2) \). Thus, the line bundle \( H^{4,0}(X \times T^2) \simeq H^{3,0}(X) \otimes H^{1,0}(T^2) \to \mathcal{C} \) is an exterior tensor product and the WP metric is a direct product. We denote an element of \( H^{3,0}_\tau(X) \) by \( \Omega_z \), and an element of \( H^{1,0}_0(T^2) \) by \( \omega_{\tau} \). We often assume that \( \omega_{\tau} = d\tau + \tau d\bar{\tau} \).

### 2.2. Variational derivatives and covariant derivatives.

The bundle \( H^{3,0}_\tau(X) \to \mathcal{M} \) is a holomorphic line bundle. Since \( H^{3,0}_\tau(X) \subset H^3(X, \mathbb{C}) \), one can view a holomorphically varying family \( \Omega_z \in H^{3,0}_\tau(X) \) as a holomorphic map \( \mathcal{M} \to H^3(X, \mathbb{C}) \) or as a holomorphic section of \( H^{3,0}_\tau(X) \). As a holomorphic vector valued function, \( \Omega_z \) may be differentiated in \( z \). If \( z_1, \ldots, z_{h^{1,1}} \) are local holomorphic coordinates, and if \( \{ \frac{\partial}{\partial z_j} \} \) are the coordinate vector fields, then \( \frac{\partial \Omega_z}{\partial z_j} \) is a well-defined element of \( H^3(X, \mathbb{C}) \).

By the Griffiths transversality theorem, see [GHJ, CO], (5.4) or [Wa1, Wa2],

\[
\frac{\partial \Omega_z}{\partial z_j} = k_j(z) \Omega_z + \chi_j,
\] (40)

where \( \chi_j \in H^{2,1}_\tau(X) \) and where \( k \in C^\infty(\mathcal{M}) \). Note that although \( \frac{\partial \Omega_z}{\partial z_j} \) is holomorphic, neither term on the right-hand side is separately holomorphic.

We define a Levi-Civita connection on the bundle \( H^{3,0} \to \mathcal{M} \) by orthogonally projecting the derivatives \( \frac{\partial \Omega_z}{\partial z_j} \) onto \( H^{3,0} \). This defines the Weil-Petersson connection \( \nabla_{WP} \) on \( H^{3,0} \to \mathcal{M} \),

\[
\nabla_{WP} : C^\infty(\mathcal{M}, \mathcal{L}) \to C^\infty(\mathcal{M}, \mathcal{L} \otimes T^*)
\]

It follows from (40) that

\[
\frac{\partial}{\partial z_j} \int_X \Omega_z \wedge \overline{\Omega_z} = k_j \int_X \Omega_z \wedge \overline{\Omega_z},
\] (41)

which by (39) implies

\[
k_j = -\frac{\partial K}{\partial z_j}.
\] (42)

Hence,

\[
\nabla_{WP} \Omega_z = -\partial K \otimes \Omega_z = \sum k_j dz_j \otimes \Omega_z
\]
is the Chern connection of the Weil-Petersson Hermitian metric.
We also define the forms
\[
\begin{aligned}
D_j \Omega_z &= \frac{\partial}{\partial z^j} \Omega + \frac{\partial K}{\partial z^j} \Omega \\
D_j D_k \Omega_z &= \left( \frac{\partial}{\partial z^j} + \frac{\partial K}{\partial z^j} \right) \left( \frac{\partial}{\partial z^k} + \frac{\partial K}{\partial z^k} \right) \Omega_z.
\end{aligned}
\]
(43)

We then have
\[
D_j \Omega_z = \frac{\partial \Omega_z}{\partial z^j} - k_j \Omega_z = \chi_j \in H^{2,1}(X). \tag{44}
\]

The operator $D_j \Omega_z$ is analogous to the second fundamental form $II(X, Y) = (\hat{\nabla}_X Y)_{\perp}$ of an embedding, i.e. it is the ‘normal’ component of the ambient derivative. It is known that the first variational derivatives span $H^{2,1}$ (see e.g. [Wa1, Wa2]. (In the physics literature, $D_\alpha$ is often described as a connection, and is often identified with $\nabla_{WP}$, but this is not quite correct as it is applied to $\Omega_z$).

The Weil-Petersson Hermitian metric $\sum G_{i\bar{j}} dz_i d\bar{z}_j$ on $\mathcal{M}$ is the curvature $(1, 1)$-form of the Hodge bundle. From (39) and (44), we have:
\[
G_{j\bar{k}} = \frac{\partial^2}{\partial z^j \partial \bar{z}^k} K(z, \bar{z}) = -\int_{\mathcal{M}} D_j \Omega_z \wedge \bar{D}_k \Omega_z \frac{1}{\int_{\mathcal{M}} \Omega_z \wedge \bar{\Omega}_z}. \tag{45}
\]

2.3. Yukawa couplings and special geometry of the moduli space. In formula (32), the density of critical points is expressed as an integral over a space $H_Z \oplus \mathbb{C}$, where $H_Z$ is a subspace of the complex symmetric matrices $\text{Sym}(h^{2,1} + 1, \mathbb{C})$ spanned by the special matrices $\xi^j$ given in (27). Their components $\mathcal{F}_{ik}^j(z)$ are known as Yukawa couplings and defined as follows: A priori, $D_k D_j \Omega_z \in H^{2,1} \oplus H^{1,2}$, and moreover its $H^{2,1}$ component vanishes (see e.g. [CO, (5.5)]). Hence we may define $\mathcal{F}_{ik}^j$ by
\[
D_k D_j \Omega_z = -\sqrt{-1} e^K \mathcal{F}_{ik}^j \frac{D_i \Omega}{\Omega} \quad (1 \leq j, k, l \leq h^{2,1}). \tag{46}
\]

See also [St1, (28)]. It is further shown in [St1, (37)] (see also [AD, (4.8)], [LS2, Theorem 3.1]) that the Riemann tensor of the Weil-Petersson metric on the moduli space $\mathcal{M}$ of Calabi-Yau three-folds is related to the Yukawa couplings by
\[
R_{ijk\ell} = G_{ij} G_{k\ell} + G_{i\ell} G_{kj} - \varepsilon^{2k} \sum_p G^{pq} \mathcal{F}_{ikp} \mathcal{F}_{j\ell q}. \tag{47}
\]

The Yukawa couplings are related to the periods of $\Omega_z$ and to the so-called prepotential of $\mathcal{M}$. We pause to recall the basic relations and to direct the reader to the relevant references.

First, we consider periods. As a basis of $H_3(X, \mathbb{R})$ we choose the symplectic basis consisting of dually paired Lagrangian subspaces of $A$-cycles $A_a$ and $B$-cycles $B_a$. The periods of $\Omega_z \in H_3^{3,0}(X)$ over the $A$-cycles
\[
\zeta^a = \int_{A_a} \Omega_z \quad (1 \leq a \leq h^{2,1} + 1 = b_3/2)
\]
define holomorphic coordinates on $L^a_\mathcal{M} = H_3^{3,0} \to \mathcal{M}$. Alternately, we can view the $\zeta^a$ as ‘special’ projective coordinates on $\mathcal{M}$. The periods of $\Omega_z$ over the $B$-cycles are then
holomorphic functions of the $\zeta^a$. The principal fact is that the image of $L^*_M$ under the period map is a complex Lagrangian submanifold of $H^3(M, \mathbb{C})$, and thus is determined by a single holomorphic function, the “prepotential” $\mathcal{F} = \mathcal{F}(\zeta^1, \ldots, \zeta^{b_3/2}) : \mathcal{L}^*_M \rightarrow \mathbb{C}$ such that

$$\int_{B_a} \Omega_z = \frac{\partial \mathcal{F}}{\partial \zeta^a}.$$ (48)

Furthermore, $\mathcal{F}$ is homogeneous of degree 2 in the periods $\zeta^a$,

$$\sum_{j=1}^{b_3/2} \zeta^a \frac{\partial \mathcal{F}}{\partial \zeta^a} = 2\mathcal{F}(z),$$

and hence may be viewed as a holomorphic section of $L^* \otimes^2 M$.

The local holomorphic 3-form $\Omega_z$ may be expressed in terms of the Poincaré duals of the symplectic basis by:

$$\Omega_z = \sum_{a=1}^{b_3/2} \left( \zeta^a \hat{A}_a - \frac{\partial \mathcal{F}}{\partial \zeta^a} \hat{B}_a \right).$$ (49)

(See [CO], (3.8).) Further, in these coordinates, the Kähler potential (39) of the Weil-Petersson metric may be written as

$$K(z, \bar{z}) = -\log i \left( \sum_{a=1}^{b_3/2} \zeta^a \hat{A}_a - \frac{\partial \mathcal{F}}{\partial \zeta^a} \right).$$

We also have:

$$\mathcal{F}_{k\bar{i}} = \frac{h^{2,1}}{2} \sum_{r=1}^{r^{2,1}} G_{r\bar{i}} \frac{\partial^3 \mathcal{F}}{\partial z^r \partial z^j \partial z^k}.$$ (50)

See [CO, (4.5)] and [St1, (64)].

In summary, we reproduce the table from [CO]:

| Derivatives of the Basis | spans |
|--------------------------|-------|
| $\Omega$                 | $H^{3,0}$     |
| $\mathcal{D}_j \Omega$  | $H^{2,1}$     |
| $\mathcal{D}_k \mathcal{D}_j \Omega$ | $-i e^K \mathcal{F}_{k\bar{j}} \mathcal{D}_{\gamma} \Omega$ |
| $\mathcal{D}_k \mathcal{D}_j \Omega$ | $G_{k\bar{j}} \mathcal{F}$ |

$$H^{12} \quad H^{03}$$ (51)
2.3.1. $\mathcal{C}$ as the moduli space of complex structures on $X \times T^2$. Above, we have reviewed the geometry of the moduli space of complex structures on the Calabi-Yau three-fold. Our configuration space $\mathcal{C} = \mathcal{M} \times \mathcal{H}$ may be viewed as (a component of) the moduli space of complex structures on $X \times T^2$. This point of view is used in [DD1], but because the $T^2$ factor plays a distinguished role we do not emphasize this identification here. Further, formula (47) needs to be modified for the moduli space of complex structures on a Calabi-Yau four-fold. In [LS2, Theorem 3.1], the Riemann tensor of the Weil-Petersson metric on the moduli space of a Calabi-Yau manifold of arbitrary dimension is shown to be

$$R_{i\bar{j}k\bar{\ell}} = G_{ij}G_{k\bar{\ell}} + G_{i\bar{\ell}}G_{k\bar{j}} - \frac{\langle D_k D_i /\Omega, \overline{D_\ell D_j /\Omega} \rangle}{\int_{\mathcal{M}} \Omega \wedge \overline{\Omega}}.$$ (52)

In the case of three-folds, the vectors $D_j /\Omega$ form an orthonormal basis for $H^{2,1}$ and one can write the inner product in the form (47).

2.4. Hodge-Riemann form and inner products. The Hodge-Riemann bilinear form on $H^3(X, \mathbb{R})$ is the intersection form $(\alpha, \beta) \mapsto \int_X \alpha \wedge \beta$. We consider the sesquilinear pairing:

$$(\alpha, \beta) \mapsto Q(\alpha, \overline{\beta}) = -\sqrt{-1} \int_X \alpha \wedge \overline{\beta}, \quad \alpha, \beta \in H^3(X, \mathbb{C}).$$ (53)

An important fact is that under the Hodge decomposition (3) for a given complex structure, the Hodge-Riemann form is definite in each summand:

$$(-1)^p Q(\alpha, \overline{\alpha}) > 0, \quad \alpha \in H^{p,3-p}(X, \mathbb{C}),$$ (54)

whose sign depends only on the parity of $p$. (See [GH, §7]. Note that our definition of $Q$ has the extra sign $-\sqrt{-1}$. The inequality (54) holds only for primitive forms, but in our case all harmonic 3-forms are primitive, since we are assuming that $H^1(M, \mathbb{C}) = 0$.) To restate (54):

**Proposition 2.1.** Let $\dim X = 3$, and let $b_1(X) = 0$. Then for each $z \in \mathcal{M}$, the Hodge-Riemann form is positive definite on $H^{2,1}_z \oplus H^{0,3}_z$ and negative definite on $H^{3,0}_z \oplus H^{1,2}_z$.

By Griffiths transversality (see (40)), for any local holomorphic frame $\Omega_z$, $D_j \Omega_z \in H^{2,1}_z$ and these elements span $H^{2,1}_z$. Also, $\overline{\Omega}_z$ spans $H^{0,3}_z$. These forms provide us with an orthonormal basis for $H^{2,1}_z \oplus H^{0,3}_z$:

**Proposition 2.2.** If $\{z_j\}$ are coordinates at $z_0$ such that $\{\partial /\partial z_j |_{z_0}\}$ are orthonormal, and if $h_{WP}(\Omega_{z_0}, \overline{\Omega}_{z_0}) = 1$, then the basis $\{D_j /\Omega_{z_0}, \overline{\Omega}(z_0)\}$ is a complex orthonormal basis of $H^{2,1}_z \oplus H^{0,3}_{z_0}$ with respect to the Hodge Riemann form $Q$.

**Remark.** Here and below, when we say that a basis of a complex vector space is complex orthonormal we mean that it is a complex basis and is orthonormal for the given inner product. By a real orthonormal basis of the same vector space we mean an orthonormal basis of the underlying real vector space.
Proof. It suffices to show that:

(i) \[ Q(\mathcal{D}_j \Omega_z, \overline{\mathcal{D}_k \Omega_z}) = -i \int_X \mathcal{D}_j \Omega_z \wedge \overline{\mathcal{D}_k \Omega_z} = G_{j\bar{k}} e^{-K}, \]

(ii) \[ Q(\mathcal{D}_j \Omega_z, \Omega_z) = -i \int_X \mathcal{D}_j \Omega_z \wedge \Omega_z = 0, \]

(iii) \[ Q(\overline{\Omega_z}, \Omega_z) = -i \int_X \overline{\Omega_z} \wedge \Omega_z = h_{WP}(\Omega_z, \Omega_z). \]

Equation (i) follows from (45), (ii) is by type considerations, and (iii) follows from (10).

Remark. In the language of complex symplectic geometry, Proposition 2.1 says that \( H^{2,1}_z \oplus H^{0,3}_z \) is a positive complex polarization of \( H^3(X, \mathbb{C}) \). Let us recall the definitions. The space \( (H^3(X, \mathbb{R}), Q) \) of real 3-cycles with its intersection form \( Q(\alpha, \beta) = -i \int_M \alpha \wedge \beta \) is a real symplectic vector space. After complexifying, we obtain the complex symplectic vector space \( (H^3(X, \mathbb{C}), Q) \). In general, if \( (V, \omega) \) is a real symplectic vector space and if \( (V_{\mathbb{C}}, \omega_{\mathbb{C}}) \) is its complexification, a complex Lagrangian subspace \( F \subset V_{\mathbb{C}} \) is called a polarization. The polarization is called real if \( F = \overline{F} \) and complex if \( F \cap \overline{F} = \{0\} \). The polarization \( F \) is called positive if \( i\omega(v, \bar{w}) \) is positive definite on \( F \).

In our setting, \( (V, \omega) = (H^3(X, \mathbb{R}), Q) \). We observe that for any complex structure \( z \) on \( X \) (as a complex manifold), the Hodge decomposition may be written in the form

\[ H^3(X, \mathbb{C}) = F \oplus \overline{F}, \quad F = H^{2,1} \oplus H^{0,3} \quad \overline{F} = H^{3,0} \oplus H^{1,2}, \]

where \( F \) is complex Lagrangian. By Proposition 2.1, this polarization is positive, i.e.

\[ Q(v, \bar{v}) > 0, \quad v \in F \setminus \{0\}. \]

3. Critical Points of Superpotentials

In this section, we assemble some basic facts about critical points and Hessians of flux superpotentials.

3.1. Flux superpotentials as holomorphic sections. As discussed in the previous section, \( \mathcal{L} \rightarrow \mathcal{C} \) is a negative line bundle. On a compact complex manifold, a negative line bundle has no holomorphic sections. However, \((\mathcal{C}, \omega_{WP})\) is a non-compact, incomplete Kähler manifold of finite Weil-Petersson volume (see [LS1] for the latter statement), and the line bundle \( \mathcal{L} \rightarrow \mathcal{C} \) has many holomorphic sections related to the periods of \( X \times T^2 \).

The sections relevant to this article are the flux superpotentials \( W_G \) of (5)–(6). \( W_G \) depends on two real fluxes \( F, H \in H^3(X, \mathbb{Z}) \), which we combine into a complex integral flux

\[ G = F + iH \in H^3(X, \mathbb{Z} \oplus \sqrt{-1}\mathbb{Z}). \]

The main reason to form this complex combination is that it relates the tadpole constraint (1) on the pair \( (F, H) \) to the Hodge-Riemann form (2). However, none of subsequent
identifications preserves this complex structure, and the reader may prefer to view \( G \) as just the pair \( G = (F, H) \in H^3(X, \mathbb{Z}) \oplus H^3(X, \mathbb{Z}) \). Alternately, we can identify \( G = F + iH \in H^3(X, \mathbb{C}) \) with the real cohomology class \( \tilde{G} := F \wedge dy - H \wedge dx \in H^4(X \times T^2, \mathbb{R}) \approx H^3(X, \mathbb{C}) \).

We shall consider the (real-linear) embedding
\[
\mathcal{W} : H^3(X, \mathbb{C}) \to H^0(\mathcal{C}, \mathcal{L}), \quad G \mapsto W_G,
\]
where \( W_G \) is given by formula (6); i.e.,
\[
(W_G(z, \tau), \Omega_z \otimes \omega_\tau) = \int_{X \times T^2} \tilde{G} \wedge \Omega_z \wedge \omega_\tau.
\]
We denote by \( S = \text{Image}(\mathcal{W}) \) the range of this map, and by
\[
S^\mathbb{Z} = \mathcal{W}(H^3(X, \mathbb{Z} \oplus i\mathbb{Z})).
\]
the lattice of sections satisfying the integrality condition. The map \( G \to W_G \) is not complex linear, so \( S \) is not a complex subspace of \( H^0(M \times \mathcal{E}, \mathcal{L}) \). Rather, it is a real subspace of dimension \( 2b_3 \) (over \( \mathbb{R} \)) and \( S^\mathbb{Z} \) is a lattice of rank \( 2b_3 \) in it. In fact \( S \approx \mathbb{R}^{2b_3} \) is totally real in \( H^0(\mathcal{C}, \mathcal{L}) \approx \mathbb{C}^{2b_3} \).

We choose local holomorphic frames \( \Omega_z \) of the Hodge bundle \( H^3, 0 \to M \) and \( \omega_\tau = dx + \tau dy \) of \( H^1, 0 \to \mathcal{E} \) and let \( \Omega_z^* \otimes \omega_\tau^* \) denote the dual co-frame of \( \mathcal{L} \). A holomorphic section of \( \mathcal{L} \) can then be expressed as \( W = f(z, \tau)\Omega_z^* \otimes \omega_\tau^* \), where \( f \in \mathcal{O}(\mathcal{C}) \) is a local holomorphic function. If \( W = W_G \) is a flux superpotential, then the corresponding function \( f_G \) is given by:
\[
f_G(z, \tau) = \int_{X \times T^2} (F \wedge dy - H \wedge dx) \wedge (\Omega_z \wedge \omega_\tau).
\]

When \( \omega_\tau = dx + \tau dy \) (on a fundamental domain in Teichmüller space), we obtain the simpler form:
\[
f_G(z, \tau) = \int_X (F + \tau H) \wedge \Omega_z.
\]

### 3.2. Critical points and Hessians of holomorphic sections

As preparation for critical points of superpotentials, we recall some basic notations and facts concerning critical points and Hessians of holomorphic sections of a general line bundle \( L \to M \) (see [DSZ1]).

Let \( (L, h) \to M \) be a holomorphic Hermitian line bundle, let \( e_L \) denote a local frame over an open set \( U \) and write a general holomorphic section as \( s = fe_L \) with \( f \in \mathcal{O}(U) \). Recall that the Chern connection \( \nabla_h \) of \( h \) is given locally as \( \nabla_h(f e_L) = (\partial f - f \partial K) \otimes e_L \), where \( K = -\log \|e_L\|^2_h \), i.e.
\[
\nabla_h = \sum_{j, l=1}^m \left( \frac{\partial f}{\partial z_l} e_L \right) dz_j \otimes e_L = \sum_{j=1}^m e_L \frac{\partial}{\partial z_j} (e^{-K} f) dz_j \otimes e_L.
\]

\[ (57) \]
The critical point equation thus reads,
\[ \frac{\partial f}{\partial z_j} - f \frac{\partial K}{\partial z_j} = 0. \]

The Hessian of a holomorphic section \( s \) of \((L, h) \to M\) at a critical point \( Z_0 \) is the tensor
\[ D \nabla s(Z_0) \in T^* \otimes T^* \otimes L, \]
where \( D \) is a connection on \( T^* \otimes L \). At a critical point \( Z_0 \), \( D \nabla s(Z_0) \) is independent of the choice of connection on \( T^* \). In a local frame and in local coordinates we have
\[
D' \nabla' s(Z_0) = \sum_{j,q} H'_{jq} dz^q \otimes dz^j \otimes e_L, \\
D'' \nabla' s(Z_0) = \sum_{j,q} H''_{jq} d\bar{z}^q \otimes dz^j \otimes e_L. \tag{58}
\]
The Hessian \( D \nabla s(Z_0) \) at a critical point thus determines the complex symmetric matrix \( H^c \) (which we call the ‘complex Hessian’):
\[
H^c := \begin{pmatrix}
H' & H'' \\
H'' & H'
\end{pmatrix} = \begin{pmatrix}
H' - f(Z_0)\Theta & 0 \\
0 & -f(Z_0)\Theta
\end{pmatrix}, \tag{59}
\]
whose components are given by
\[
H'_{jq} = \left( \frac{\partial}{\partial z^j} - \frac{\partial K}{\partial z^j} \right) \left( \frac{\partial}{\partial z^q} - \frac{\partial K}{\partial z^q} \right) f(Z_0), \tag{60}
\]
\[
H''_{jq} = -f \left. \frac{\partial^2 K}{\partial z^j \partial \bar{z}^q} \right|_{Z_0} = -f(Z_0)\Theta_{jq}, \quad \Theta_{h}(Z_0) = \sum_{j,q} \Theta_{jq} dz^j \wedge d\bar{z}^q. \tag{61}
\]

3.3. Supersymmetric critical points and the Hodge decomposition. We now specialize to the critical point equations for flux superpotentials \( W_G(z, \tau) \). An important observation that is now standard in the physics literature is that the complex moduli \((z, \tau)\) at which a flux superpotential \( W_G(z, \tau) \) satisfies \( \nabla W_G = 0 \) are characterized by the following special Hodge decomposition of \( H^3(X, \mathbb{C}) \) at \( z \) (see [AD], (3.5)–(3.8)).

A local holomorphic frame for the Hodge bundle \( \mathcal{L} \to C \) is \( e_\mathcal{L} = \Omega_z^* \otimes \omega_\tau^* \), where \( \Omega_z^* \) is dual to the local frame \( \Omega_z \) of the Hodge line bundle \( H^{3,0} \to \mathcal{M} \) and \( \omega_\tau^* \) is dual to the local frame \( \omega_\tau = dx + \tau dy \) of \( H^{1,0} \to \mathcal{E} \). We let \( K(Z) = K_X(z) + K_{T^2}(\tau) \) be the Kähler potential for the local frame \( \Omega_z \otimes \omega_\tau \) of the (positive) Hodge bundle \( \mathcal{L}^* \). We then have
\[
|e_\mathcal{L}(Z)|^2_n = |\Omega_z \otimes \omega_\tau|^{-2}_{hwp} = e^{K(Z)} = e^{K_X(z)} e^{K_{T^2}(\tau)}. \tag{62}
\]

Hence, the Weil-Petersson Kähler potential on \( \mathcal{C} \) is
\[
K(Z) = -\log \int_X \Omega_z \wedge \bar{\Omega}_z - \log(\bar{\tau} - \tau). \]
In particular, the $\tau$-covariant derivative on $\mathcal{L}$ is given in the local frame $e_\mathcal{L}$ by

$$\nabla_\tau = \frac{\partial}{\partial \tau} + \frac{1}{\tau - \bar{\tau}}. \quad (63)$$

Hence with $W_G = f_G e_\mathcal{L}$, we have

$$\nabla_\tau f_G = \int_X \left[ H + \frac{1}{\tau - \bar{\tau}} (F + \tau H) \right] \wedge \Omega_z$$

$$= \frac{1}{\tau - \bar{\tau}} \int_X (F + \bar{\tau} H) \wedge \Omega_z. \quad (64)$$

To compute the $z$-derivatives, we see from §2.2 and (56)–(57) that

$$\nabla_z f_G = \left( \frac{\partial f_G}{\partial z} + \frac{\partial K}{\partial z} f_G \right) (z, \tau) = \int_X (F + \tau H) \wedge \left( \frac{\partial \Omega_z}{\partial z} + \frac{\partial K}{\partial z} \Omega_z \right)$$

$$= \int_X (F + \tau H) \wedge \chi_j = 0, \quad (65)$$

for $1 \leq j \leq h^{2,1}$. Thus, the supersymmetric critical point equations for the flux superpotential $W_G$ read:

$$\begin{cases} \int_X (F + \tau H) \wedge D_j \Omega_z = 0 & (1 \leq j \leq h^{2,1}) \\ \int_X (F + \tau H) \wedge \overline{\Omega}_z = 0. \end{cases} \quad (66)$$

As in (16), we denote by $S_Z (Z = (z, \tau))$ the space of superpotentials $W_G$ with $\nabla W_G (Z) = 0$. Although the equation is complex linear on $H^0(\mathcal{C}, \mathcal{L})$, $S$ is not a complex subspace of $H^0(\mathcal{C}, \mathcal{L})$, so $S_Z$ is a real but not complex vector space. Put another way, for each $Z = (z, \tau)$, the critical point equation determines a real subspace

$$H^3_Z (X, \mathbb{C}) = \mathcal{W}^{-1} (S_Z) = \{ F + i H, \ F, H \in H^3 (X, \mathbb{R}), \ \ (66) \text{ is true} \}. \quad (67)$$

The critical point equations (66) put $b_3 = 2(h^{2,1} + 1)$ independent real linear conditions on $2b_3$ real unknowns $(F, H)$.

**Proposition 3.1 ([AD, DD1]).** Let $G = F + i H$ with $F, H \in H^3 (X, \mathbb{R})$, and let $\langle W_G (z, \tau), \Omega_z, \omega_z \rangle = \int_X (F + \tau H) \wedge \Omega_z$ be the associated superpotential. If $\nabla_{z, \tau} W_G (z, \tau) = 0$, then $(F + \tau H) \in H^{2,1}_z \oplus H^{0,3}_z$. Moreover, the map

$$I_\tau : H^3 (X, \mathbb{C}) \to H^3 (X, \mathbb{C}), \quad I_\tau (F + i H) = F + \tau H$$

restricts to give real linear isomorphisms

$$I_{z, \tau} : H^3_{z, \tau} \to H^{2,1}_z (X) \oplus H^{0,3}_z (X),$$

of real vector spaces.

**Proof.** We first prove that $(F + i H) \mapsto F + \tau H$ takes $H^3_Z \mapsto H^{2,1}_z \oplus H^{0,3}_z$. Suppose that $\nabla W_G = 0$. Since the $\chi_j (z)$ span $H^{2,1}_z$, we conclude from the first equation of (66) that $(F + \tau H)_1 = 0$; by the second equation, we also have $(F + \tau H)_0 = 0$. Thus $F + \tau H \in H^{2,1}_z \oplus H^{0,3}_z$.

Since $I_{z, \tau}$ is injective and since $\dim \mathbb{R} H^3_{z, \tau} = \dim \mathbb{R} H^{2,1}_z \oplus H^{0,3}_z = b_3$, it is clearly an isomorphism. $\square$
3.4. The map \((z, \tau) \to H^3_{z,\tau}\). As \((z, \tau)\) varies over \(\mathcal{C}\), how do the spaces \(H^3_{z,\tau}\) move in \(H^3(X, \mathbb{C})\)? This question is important in relating the pure lattice point problem in \(H^3(X, \mathbb{C})\) to the vacuum distribution problem in \(\mathcal{C}\). It depends on the geometry of the diagram

\[
\mathcal{I} \subset \mathcal{C} \times H^3(X, \mathbb{C})
\]

where \(\mathcal{I} = \{(z, \tau, F, H) : F + iH \in H^3_{(z,\tau)}(X)\}\), which is a replica of (15) in which \(S\) is replaced by \(H^3(X, \mathbb{C})\).

To answer this question, we first note that for each \((z, \tau) \in \mathcal{C}\), the real-linear map \(H^3_{z,\tau} \to H^3(X, \mathbb{R}), F + iH \mapsto H\) is bijective. Injectivity follows by noting that

\[
F \in H^3_{z,\tau} \implies F \in H^2_{z,1} \oplus H^0_{z,3} \implies F = \tilde{F} \in H^{1,2}_{z} \oplus H^3_{z,0} \implies F = 0.
\]

Since both spaces have dimension \(b_3\), bijectivity follows. Thus there is a real linear isomorphism \(\iota_{z,\tau}(H) = F(z, \tau, H) + iH\).

To describe \(F(z, \tau, H)\), we form the \(z\)-dependent basis

\[
\{\text{Re } D_1 \Omega_z, \ldots, \text{Re } D_{h^2,1} \Omega_z, \text{Re } \Omega_z, \text{Im } D_1 \Omega_z, \ldots, \text{Im } D_{h^2,1} \Omega_z, -\text{Im } \Omega_z\}
\]

of \(H^3(X, \mathbb{R})\). We then have

\[
F(z, \tau, H) = J_{\tau} H,
\]

where \(J_{\tau}\) is given by the block matrix

\[
J_{\tau} = \begin{pmatrix}
\text{Re } \tau I_m & -\text{Im } \tau I_m \\
\text{Im } \tau I_m & \text{Re } \tau I_m
\end{pmatrix},
\]

with respect to the basis (69).

This yields the following proposition:

**Proposition 3.2.** The mapping \((z, \tau, H) \mapsto (z, \tau, \iota_{z,\tau}(H))\) gives an isomorphism \(\mathcal{C} \times H^3(X, \mathbb{R}) \cong \mathcal{I}\).

An important consequence is:

**Proposition 3.3.** For any open subset \(U \subset \mathcal{C}\), the cone \(\bigcup_{(z,\tau) \in U} H^3_{(z,\tau)}(X) \setminus \{0\}\) is open in \(H^3(X, \mathbb{C}) \setminus \{0\}\).
Proof. We must show that
\[ \pi \left[ I \cap \{ U \times H^3(X, \mathbb{C}) \} \right] \setminus \{ 0 \} \]
is open. By Proposition 3.2, it suffices to show that the image of the map
\[ \iota : U \times \{ H^3(X, \mathbb{R}) \setminus \{ 0 \} \} \to H^3(X, \mathbb{C}), \quad \iota(z, \tau, H) = \iota_{z, \tau}(H) = F(z, \tau, H) + iH \]
is open. We fix \((z_0, \tau_0, H_0)\) and consider the derivative \(D\iota|_{z_0,\tau_0,H_0}\) on \(T_{z_0,\tau_0}C \times H^3(X, \mathbb{R}) \setminus \{ 0 \} \). Since the linear map \(\iota_{z,\tau}\) is bijective, if we vary \(H\), we get all of \(H^3_{z,\tau}\), so the issue is to prove that we obtain the complementary space by taking variations in \(\tau, z\).

First, \(H^3_{z,\tau} = I_\tau^{-1}(H^2_{z,1} \oplus H^0_{z,3})\). The \(z\) variations of \(H^2_{z,1} \oplus H^0_{z,3}\) span this space plus \(H^1_{z,2}\). By (69)–(71), variations in Re \(\tau\), resp. Im \(\tau\), produce Re \(\Omega_{z}\), Im \(\Omega_{z}\) and hence \(H^3_{z,0} = \text{span}(\Omega_{z})\) is also in the image. \(\square\)

Remark. We could also ask what kind of set is swept out in \(\bigcup_{z \in U} H^2_{z,1} \oplus H^0_{z,3}\) as \(z\) ranges over an open set \(U \subset M\). Since \(\dim \mathbb{C} U = h^{2,1}\), the image of this map is a real codimension two submanifold.

3.5. Inner product on \(S_Z\). In Theorem 1.4, combined with (32), we have expressed \(N_\psi(L)\) in terms of a Gaussian type ensemble of holomorphic sections in \(S_Z\). We now clarify the inner product \(Q_z\) (cf. (18)), Gaussian measure and Szegö kernel on this space.

**Proposition 3.4.** The Hodge-Riemann Hermitian inner product on \(H^3(X, \mathbb{C})\) restricts for each \(Z = (z, \tau)\) to define a complex valued inner product on \(H^3_Z\) which satisfies \(Q_Z[G] > 0\) for all \(G \neq 0\). Moreover, the map \(I_\tau : H^3_Z \to H^2_{z,1} \oplus H^0_{z,3}\) satisfies \(Q[I_\tau G] = \text{Im} \tau \cdot Q[G]\).

**Proof.** It follows by Proposition 2.1 that the symmetric bilinear form
\[ Q[F + \tau H] = i^3 \int_X (F + \tau H) \wedge (\bar{F} + \bar{\tau} \bar{H}) = \text{Im} \tau \cdot Q[F + iH] \]
on \(H^3_{z,\tau}(X, \mathbb{C})\) in (67) is positive definite. \(\square\)

Recall that we have the real-linear isomorphisms
\[ H^3(X, \mathbb{C}) \xrightarrow{\mathcal{W}} S \subseteq H^0(C, L) \]
\[ I_\tau \downarrow \]
\[ H^3(X, \mathbb{C}) \]
where \(I_\tau(F + iH) = F + \tau H\). Restricting (73) to fluxes with a critical point at \(Z = (z, \tau)\), we have isomorphisms
\[ H^3_Z \xrightarrow{\mathcal{W}} S_Z 
I_\tau \downarrow
H^2_{z,1} \oplus H^0_{z,3} \]

(74)
We let $\tilde{Q}$ denote the Hermitian inner product on $H^{2,1}_\mathbb{C} \oplus H^{0,3}_\mathbb{C}$ transported from $(H^3_\mathbb{C}, Q)$ by $I_\tau$; i.e.,

$$\tilde{Q}[C] = Q\left[I^{-1}_\tau(C)\right], \quad C \in H^{2,1}_\mathbb{C} \oplus H^{0,3}_\mathbb{C}. \quad (75)$$

Hence by (72), we have:

$$Q[C] = \left(\text{Im } \tau\right) \tilde{Q}[C]. \quad (76)$$

As in (18), we define the inner product $Q_Z$ on $S_Z$ by $Q_Z[W_G] = Q_Z[G]$. We also denote by $\Pi_Z$ the Szegő kernel of $S_Z$, i.e. the orthogonal projection onto this subspace with respect to $Q_Z$. An explicit formula for $\Pi_Z$ is given in [AD] (4.4) in terms of covariant derivatives of the Kähler potential. As with any inner product space, the inner product $Q_Z$ induces a Gaussian measure on $S_Z$ (see [DSZ1]). But we emphasize that we are counting flux vacua and not studying probabilities with respect to a normalized Gaussian measure.

4. Counting Critical Points: Proof of Proposition 1.3

We now prove the first result on counting critical points of flux superpotentials $W_G$ where $G$ satisfies the tadpole constraint (1). Before starting the proof, we review the geometry of the lattice point problem and the critical point problem.

We wish to count vacua in a region of moduli space as $G$ varies over fluxes satisfying the tadpole constraint. Equivalently, we count inequivalent vacua in Teichmüller space. That is, $\Gamma$ acts on the pairs $(W, Z)$ of superpotentials and moduli by

$$\gamma \cdot (G, Z) = (\varphi(\gamma) \cdot G, \gamma \cdot Z).$$

Therefore $\Gamma$ acts on the incidence relation (14). We only wish to count critical points modulo the action of $\Gamma$. To do this, there are two choices: we could break the symmetry by fixing a fundamental domain $D_\Gamma \subset \mathcal{C}$ for $\Gamma$ in $\mathcal{C}$, i.e. only count critical points in a fundamental domain. Or we could fix a fundamental domain for $\varphi(\Gamma)$ in $H^3(X, \mathbb{C})$ and count all critical points of these special flux superpotentials. When we do not know the group $\varphi(\Gamma)$ precisely, it seems simpler to take the first option and that is what we do in Proposition 1.3 and Theorem 1.4. We note that the number of critical points of $W_G$ in Teichmüller space equals the number of critical points of the $\Gamma$-orbit of $W_G$ in $\mathcal{C}$.

The level sets $Q[G] = C$ for $C > 0$ are hyperboloids contained in $\{G : Q[G] > 0\}$ and thus the tadpole constraint defines a hyperbolic shell in $\{G : Q[G] > 0\}$. The critical point equation $\nabla W_G(Z) = 0$ is homogeneous of degree 1 in $G$. Hence, summing a homogeneous function over $G \in \{G : Q[G] > 0\}$ with $Q[G] \leq L$ may be viewed as summing a function on the hyperboloid $Q[G] = 1$ over the radial projections of the lattice points $G$ in the shell $Q[G] \leq L$. The number which project over a compact subset of $Q[G] = 1$ is finite.

4.1. Approximating the sum by an integral. Our main argument in the proof of Proposition 1.3 is the following lemma:

**Lemma 4.1.** Let $\psi = \chi_K$, where $K \subset I$ is as in Proposition 1.3. Then

$$N_\psi(L) = L^{b_3} \left[ \int_S (C_W, \psi) \chi_Q(W) dW + O \left(L^{-1/2}\right) \right].$$
Proof. We consider the integer-valued function
\[ f(W) = \langle C_W, \psi \rangle = \sum_{\{Z : \nabla W(Z) = 0\}} \psi(Z, W) = \#\{Z \in C : (Z, W) \in K\}. \]

We note that the characteristic function of the set \( \{0 \leq Q[W] \leq L\} \) is \( \chi_Q(W/\sqrt{L}) \). Using our symplectic basis to identify \( H^3(X, \mathbb{Z} \oplus \sqrt{-1} \mathbb{Z}) \) with \( \mathbb{Z}^{2b_3} \), we have
\[ N_{\psi}(L) = \sum_{N \in \mathbb{Z}^{2b_3}} f(N) \chi_Q(N/\sqrt{L}) \]
\[ = \sum_{N \in \mathbb{Z}^{2b_3}} f(N/\sqrt{L}) \chi_Q(N/\sqrt{L}) = \sum_{N \in \mathbb{Z}^{2b_3}} g(N/\sqrt{L}), \]
where
\[ g = f \chi_Q. \]

We note that \( f \) is constant on each connected component of \( S \setminus [D \cup \pi(\partial K)] \). Since the number of these components is finite, \( f \) is bounded. We let
\[ A_{\psi} = \sup_{Z \in \rho(\text{Supp } \psi)} \| Q^{-1}_Z \| < +\infty. \]

Then
\[ \inf \left\{ Q[W] : W \in \bigcup_{Z \in \rho(\text{Supp } \psi)} S(S_Z) \right\} = 1/A_{\psi} > 0. \]

Now let
\[ Q_0 := \{ W : Q[W] \leq 1, |W| \leq A_{\psi} \} \supset \text{Supp } g. \]

Approximating sums by integrals, we have
\[ L^{-b_3} N_{\psi}(L) = L^{-b_3} \sum_{N \in \mathbb{Z}^{2b_3}} g(N/\sqrt{L}) = \int_{\mathbb{R}^{2b_3}} g(W) \, dW + \sum_{N \in \mathbb{Z}^{2b_3}} E_{N,L}, \]
where
\[ E_{N,L} = \int_{\mathcal{R}_{N,L}} [g(N/\sqrt{L}) - g(W)] \, dW, \]
\[ \mathcal{R}_{N,L} = \{ W = (W_1, \ldots, W_{2b_3}) \in \mathbb{R}^{2b_3} : N_j < W_j < N_j + 1/\sqrt{L} \}. \]

Let
\[ B = Q_0 \cap [\partial Q \cup D \cup \pi(\partial K)]. \]

Since \( g \) is locally constant on \( S \setminus B \), the error \( E_{N,L} \) vanishes whenever \( \mathcal{R}_{N,L} \cap B = \emptyset \). Hence
\[ \sum_{N \in \mathbb{Z}^{2b_3}} E_{N,L} \leq (\sup f)L^{-b_3} \left[ \#\{N : \mathcal{R}_{N,L} \cap B \neq \emptyset\} \right] \]
\[ = L^{-b_3} O\left(\sqrt{L}^{2b_3-1}\right) = O(L^{-1/2}). \]
\[ \square \]
4.1.1. The index density. By applying precisely the same argument for $\text{Ind}_\psi(L)$, we obtain

**Lemma 4.2.** Let $\psi = \chi_K$, where $K \subset \mathcal{I}$ is as in Proposition 1.3. Then

$$\text{Ind}_\psi(L) = L^{b_3} \left[ \int_{\{Q[W] \leq 1\}} \langle \text{Ind}_W, \psi \rangle dW + O \left( L^{-1/2} \right) \right].$$

4.1.2. Non-clustering of critical points. Before concluding the proof of Proposition 1.3, we briefly consider the question of whether there exist real hypersurfaces $\Gamma \subset \mathcal{C}$ with the property that $\sim \sqrt{L^{2b_3-1}}$ critical points of norm $\leq L$ cluster within a $1/L$ tube around $\Gamma$. A domain in $\mathcal{C}$ whose boundary contained a piece of $\Gamma$ would attain the remainder estimate in Proposition 1.3.

Since the number of critical points corresponding to $G \in H^3(X, \mathbb{Z} \oplus \sqrt{-1} \mathbb{Z})$ is bounded, such clustering of critical points could only occur if a sublattice of rank $2b_3 - 1$ clustered around the hypersurface

$$\bigcup_{(z, \tau) \in \Gamma} H^2_{z, \tau} \subset H^3(X, \mathbb{C}). \quad (80)$$

There do exist real hypersurfaces in $H^3(X, \mathbb{C})$ for which such exceptional clustering occurs, namely hyperplanes containing a sublattice of rank $2b_3 - 1$. We refer to such a hyperplane as a rational hyperplane $L$. For instance, any pair of integral cycles $\gamma_1, \gamma_2$ defines a rational hyperplane

$$L = L_{\gamma_1, \gamma_2} = \{ G = F + iH \in H^3(X, \mathbb{C}) : \ell(F + iG) := \int_{\gamma_1} F + \int_{\gamma_2} H = 0 \}.$$

As mentioned in the introduction, projections of the lattice points $H^3(X, \mathbb{Z} \oplus \sqrt{-1} \mathbb{Z})$ to $\partial Q$ concentrate to sub-leading order $\sqrt{L^{2b_3-1}}$ around the hypersurface of $\partial Q$ obtained by intersecting it with a rational hyperplane.

However, rational hyperplanes never have the form (80). Indeed, under the correspondence $\rho \circ \pi^*$ defined by the diagram (68), the image of a hyperplane always covers a region and not a hypersurface of $\mathcal{C}$. That is,

$$\dim(L \cap H^3_{z, \tau}) > 1 \ \forall (z, \tau) \in \mathcal{C}.$$ 

Indeed, under the identification $H^3_{z, \tau} \simeq H^3(X, \mathbb{R})$, $L|_{H^3_{z, \tau}}$ becomes the real linear functional $L(H) = \int_{\gamma_1} F(z, \tau, H) + \int_{\gamma_2} H$ on $H^3(X, \mathbb{R})$. Here, we use that $F(z, \tau, H)$ is linear in $H$. Hence, $\dim(L \cap H^3_{z, \tau}) \geq b_3 - 1$ for any $(z, \tau)$.

As will be studied in [Ze1], clustering to order $\sqrt{L^{2b_3-1}}$ can only occur if the second fundamental form of (80) is completely degenerate. Hence the fact that rational hyperplanes never have this form is strong evidence that there are no smooth hypersurfaces $\Gamma \subset \mathcal{C}$ for which lattice points cluster to subleading order around (80).
4.2. Hessians and density of critical points. The final step in the proof of Proposition 1.3 is to change the order of integration over $C$ and over $S_Z$:

**Lemma 4.3.** We have:

$$\int_{\{Q[\psi] \leq 1\}} \langle C_W, \psi \rangle dW = \int_C \int_{S_Z} \psi(Z, W) \left| \det H^{cW}(Z) \chi_Q(W) \right| dW d\text{Vol}_{WP}(Z).$$

Combining the formulas in Lemmas 4.1 and 4.3, we obtain the formula of Proposition 1.3.

The proof of Lemma 4.3 is in two parts. The first is an elementary exercise in changing variables in an integral, which we accomplish below by relating both sides to pushforwards from the incidence variety in the diagram (15). The second part involves special geometry, and is given in the next section.

We may interpret the integral

$$\int_{\{Q[\psi] \leq 1\}} \langle C_W, \psi \rangle dW$$

as an integral over $I$ as follows. Implicitly, it defines a measure $d\mu_I$ so that

$$\int_I \psi(Z, W) d\mu_I = \int_{\{Q[\psi] \leq 1\}} \langle C_W, \psi \rangle dW.$$

(81)

The measure $d\mu_I$ may be expressed in terms of the Leray measure $dL_I$ defined by a measure $d\nu$ on $S$ and the ‘evaluation map’ $\epsilon: (Z, W) \in C \times S \to \nabla W(Z)$.

The Leray form is the quotient $dL_I := \frac{dV_{WP} \times dv}{d\epsilon}$, i.e. the unique form satisfying

$$dL_I \times d\epsilon = dV_{WP} \times dv.$$

This measure is often written $\delta(\nabla W(Z))dWdV(Z)$ in the physics literature.

As suggested by the physics formula, $d\mu_I = \nabla s(Z)^* \delta_0$. However, this formula is somewhat ambiguous. If we regard $s$ as fixed, then it is simply the pullback of $\delta_0$ under $Z \to \nabla s(Z)$. It is then well-known that

$$\nabla s^* \delta_0 = \sum_{Z: \nabla s(Z) = 0} \frac{\delta_Z}{|\det H^c s(Z)|}.$$

(82)

However, when interchanging the order of integration, we really wish to think of it as a function of $s$ for fixed $Z$. So we now have a function $\epsilon_Z(s) = \nabla s(Z)$ which may be viewed as

$$\epsilon_Z: S \to \mathbb{C}^m \equiv \mathbb{R}^{b_3},$$

where $m = h^{2,1} + 1 = \frac{1}{2}b_3$. So now the zero set $\epsilon_Z^{-1}(0)$ is the subspace $S_Z$ rather than the discrete set $\text{Crit}(s)$.

To simplify the notation, we now consider the general situation where we have a real $n$-dimensional manifold $M$ and a space $S$ of functions $F: M \to \mathbb{R}^n$. In our case, $F = \nabla s$ and $M$ is a coordinate neighborhood in $C$ where $M$ has local coordinates $(x_1, \ldots, x_n)$ and $\mathcal{L}$ has a local frame. Suppose that $0$ is a regular value of $F$, so that $F$ is
a local diffeomorphism in a neighborhood $U$ of any point $x_0$ of $F^{-1}(0)$. Let $h = F^{-1}ig|_U$ in a neighborhood of 0. Then for $\varphi$ supported in a neighborhood of $x_0$, put

$$
\langle F^*\delta_0, \varphi \rangle = \langle \delta_0, \varphi(h(y))| \det dh(y) \rangle.
$$

Let $\dim \mathbb{R} S = d \geq n$. In our case, $d = 2b_3 > n = b_3$, so we introduce a supplementary linear map: for a point $u \in U \subset M$, $S_u$ is the kernel of $\varepsilon_u$, and we supplement $\varepsilon_u$ with the projection $\Pi_u : S \to S_u$. Then,

$$(\varepsilon_u, \Pi_u) : S \to \mathbb{R}^n \oplus S_u$$

is a linear isomorphism. Hence it equals its derivative, so

$$
\langle \varepsilon_u^*\delta_0, \varphi \rangle = \langle \delta_0, \varphi((\varepsilon_u, \Pi_u)^{-1})| \det(\varepsilon_u, \Pi_u)^{-1} \rangle.
$$

Now, $S$ is equipped with an inner product, which induces an inner product on $\mathbb{R}^n \oplus S_u$. We choose an orthonormal basis $\{S_1, \ldots, S_n\}$ of $S_u^\perp$, and $\{S_{n+1}, \ldots, S_d\}$ for $S_u$. Since $\Pi_u : S_u \to S_u$ is the identity, $(\varepsilon_u, \Pi_u)$ has a block diagonal matrix relative to the bases of $S = S_u^\perp \oplus S_u$ and $\mathbb{R}^n \oplus S_u$, with the identity in the $S_u$-$S_u$ block. Hence, $\det(\varepsilon_u, \Pi_u) = \det(\varepsilon_u|_{S_u^\perp})$, where the determinant is with respect to these bases.

The general case of formula (81) states that

$$
d\mu_I = |\det DW(u)| \times \frac{\chi_Q du \times dW}{d\varepsilon}.
$$

We then compute the $I$ integral as an iterated integral using the other singular fibration $\pi$, i.e. by first integrating over the fibers $S_u$:

$$
\int I \psi(u) d\mu_I = \int_U \int_{S_u} \frac{\psi(u)}{|\det(\varepsilon_u|_{S_u^\perp})|} \chi_{Q_u}(W)| \det DW(u)| dW du.
$$

Returning to our case where $F = \nabla s$, (84) becomes

$$
\int I \psi(Z) d\mu_I = \int_{\mathcal{C}} \int_{S_Z} \frac{\psi(Z, W)}{|\det(\varepsilon_Z|_{S_Z^\perp})|} |\det H^*W(Z)| \chi_{Q_Z}(W) dW d\text{Vol}_{WP}(Z).
$$

4.3. Completion of the proof of Lemma 4.3. To complete the proof of the lemma, we need to show that $|\det(\varepsilon_Z|_{S_Z^\perp})| = 1$ with respect to normal coordinates and an adapted frame at $Z_0 = (z_0, \tau_0) \in \tilde{M}$.

Recalling (73)–(74), we write

$$
\mathcal{S}_{Z_0} \setminus = I_{\tau} \circ \mathcal{W}^{-1}(\mathcal{S}_{Z_0}^\perp) = H_{\tau}^{3,0} \oplus H_{\tau}^{1,2}.
$$

A complex orthonormal basis for $\mathcal{S}_{Z_0}^\perp$ relative to $Q$ is $\{\tilde{x}_0, \tilde{x}_1, \ldots, \tilde{x}_{h^{2,1}}\}$, where $\tilde{x}_0 = \Omega_{z_0}$. A basis (over $\mathbb{R}$) for $S_{Z_0}^\perp$ is

$$
\mathcal{U}_j := \mathcal{W} \circ I_{\tau}^{-1}(\tilde{x}_j), \quad \mathcal{V}_j := \mathcal{W} \circ I_{\tau}^{-1}(\sqrt{-1} \tilde{x}_j), \quad 0 \leq j \leq h^{2,1}.
$$
The basis \(\{U_j, V_j\}\) is orthogonal with respect to \(Q_{Z_0}\), but not orthonormal. By (76),

\[
\begin{align*}
Q[U_j] &= \tilde{Q}[\tilde{x}_j] = \frac{1}{\text{Im } \tau}, \\
Q[V_j] &= \tilde{Q} \left[\sqrt{-1} \tilde{x}_j\right] = \frac{1}{\text{Im } \tau}.
\end{align*}
\]

(86)

To compute \(\det(\epsilon_{Z_0})\), we let \((z_1, \ldots, z_{h^{2,1}})\) be normal coordinates about \(z_0 \in M\), and we let \(\nabla_j f\) be given by

\[
\nabla_{\partial/\partial z_j} (f e_L) = (\nabla_j f) \otimes e_L,
\]

for \(1 \leq j \leq h^{2,1}\). We find it convenient to use the coordinate \(\tau \in \mathcal{E}\), although it is not normal, and we use the normalized covariant derivative

\[
\nabla_0 := (\text{Im } \tau) \nabla_\tau.
\]

(87)

Now we write

\[
\begin{align*}
U_j &= f_j(z) \Omega_z^* \otimes \omega_\tau^*, \\
V_j &= g_j(z) \Omega_z^* \otimes \omega_\tau^*,
\end{align*}
\]

where the local frame \(\Omega_z\) is normal at \(z_0\), and \(\omega_\tau = dx + \tau dy\). Note that the Weil-Petersson norm \(|\omega_\tau^*|\) is given by

\[
|\omega_\tau^*| = |dx + \tau dy|^{-1} = \frac{1}{(\text{Im } \tau)^{1/2}}.
\]

(88)

Taking into account (86)–(88), the \(\text{Im } \tau\) factors cancel out, and we obtain

\[
\det(\epsilon_{Z_0}) = \det \begin{pmatrix}
\text{Re } \nabla_j f_k & \text{Re } \nabla_j g_k \\
\text{Im } \nabla_j f_k & \text{Im } \nabla_j g_k
\end{pmatrix}_{j,k = 1}^{h^{2,1}}.
\]

We now evaluate the entries of the matrix. By Proposition 2.2, we have

\[
\nabla_k f_j(Z) = \int_X \frac{D_j \Omega_{z_0} \wedge D_k \Omega_z}{\Omega_{z_0} \wedge \Omega_z}, \quad \nabla_k g_j(Z) = \int_X i \frac{D_j \Omega_{z_0} \wedge D_k \Omega_z}{\Omega_{z_0} \wedge \Omega_z},
\]

and hence

\[
\nabla_j f_k(Z_0) = -i \delta_{jk}, \quad \nabla_j g_k(Z_0) = \delta_{jk}, \quad \text{for } j, k \geq 1.
\]

Also

\[
\nabla_k f_0 = \int_X \Omega_{z_0} \wedge [D_k \Omega_{z_0} - (\partial K/\partial z_j) \Omega_{z_0}] = 0, \quad \nabla_k g_0 = i \nabla_k f_0 = 0 \quad \text{for } k \geq 1.
\]

By (64), we have

\[
\begin{align*}
\nabla_0(f_j) &= (\text{Im } \tau) \nabla_\tau(f_j) = \int_X D_j \Omega_{z_0} \wedge \Omega_{z_0} = 0, \\
\nabla_0(g_j) &= -i \int_X \Omega_{z_0} \wedge \Omega_{z_0} = 0, \quad j \geq 1,
\end{align*}
\]
\[ \nabla_0(f_0) = \int_X \Omega_{z_0} \wedge \Omega_{z_0} = i, \quad \nabla_0(g_0) = \int_X i \Omega_{z_0} \wedge \Omega_{z_0} = 1. \]

Therefore,

\[ |\det(\varepsilon_{Z_0}|S_{Z_0}^\perp)| = |\det \begin{pmatrix} 0 & 1 \\ D(1, -1, \ldots, -1) & 0 \end{pmatrix}| = 1. \]

\section*{5. Proof of Theorem 1.4}

In this section we prove Theorem 1.4, which is a combination of an equidistribution theorem for radial projections of lattice points and an equidistribution theorem for critical points.

\subsection*{5.1. A local van der Corput theorem.}

We first illustrate the method of proof of Theorem 1.4 by providing a van der Corput type asymptotic estimate for the radial distribution of lattice points (Theorem 5.1). The estimate has much in common with the classical van der Corput estimate of Hlawka, Randol and others on lattice points in dilates of smooth convex sets (see for example, [Ra, Hl]), and we adapt the proof of the classical estimate to obtain our asymptotic equidistribution theorem.

Let \( Q \subset \mathbb{R}^n \) (\( n \geq 2 \)) be a bounded, smooth, strictly convex set with \( 0 \in Q^\circ \). Let \( |X|_Q \) denote the norm of \( X \in \mathbb{R}^n \) given by

\[ Q = \{ X \in \mathbb{R}^n : |X|_Q < 1 \}. \] (89)

To measure the equidistribution of projections of lattice points, we consider the sums

\[ S_f(t) = \sum_{k \in \mathbb{Z}^n \cap tQ \setminus \{0\}} f \left( \frac{k}{|k|_Q} \right), \quad \text{with} \quad f \in C^\infty(\partial Q), \quad t > 0. \]

We extend \( f \) to \( \mathbb{R}^n \) as a homogeneous function of degree 0, so that \( f(k) = f \left( \frac{k}{|k|_Q} \right) \).

Our purpose is to obtain the following asymptotics of \( S_f(t) \):

\textbf{Theorem 5.1.}

\[ S_f(t) = t^n \int_Q f \ dX + O \left( t^{n-\frac{2n}{n+1}} \right), \quad t \to \infty. \]

From this it is simple to obtain asymptotics of \( S_f(t) \) when \( f \in C^\infty(\partial Q) \) is extended as a homogeneous function of any degree \( \alpha \) to \( \mathbb{R}^n \):

\textbf{Corollary 5.2.}

Let \( f \in C^\infty(\mathbb{R}^n \setminus \{0\}) \) be homogeneous of degree \( \alpha > 0 \), and let

\[ S_f(t) = \sum_{k \in \mathbb{Z}^n \cap tQ} f(k), \quad t > 0. \]

Then

\[ S_f(t) = t^{n+\alpha} \int_Q f \ dX + O \left( t^{n-\frac{2n}{n+1}+\alpha} \right), \quad t \to \infty. \]
5.1.1. Littlewood-Paley. To deal with the singularity of \( f \) at \( x = 0 \) we use a dyadic Littlewood-Paley decomposition in the radial direction. Let \( \eta \in C_0^\infty \) with \( \eta(r) = 1 \) for \( r \leq 1 \) and with \( \eta(r) = 0 \) for \( r \geq 2 \). We then define
\[
\rho \in C_0^\infty(\mathbb{R}), \quad \rho(r) = \eta(r) - \eta(2r).
\]
Then \( \rho(r) \) is supported in the shell \( 1/2 \leq r \leq 2 \), hence \( \rho(2^j r) \) is supported in the shell \( 2^{-j-1} \leq r \leq 2^{-j+1} \). We then have:
\[
\eta(r) = \sum_{j=0}^{\infty} \rho(2^j r), \quad (r \neq 0).
\]
Indeed,
\[
\sum_{j=0}^{J} \rho(2^j r) = \eta(r) - \eta(2^J r) \to \eta(r)
\]
by the assumption that \( \eta \in C_0^\infty \).

We then write
\[
S_f(t) = \sum_{k \in \mathbb{Z}^n} f(k) \chi_{[0,1]} \left( \frac{|k|_Q}{t} \right) = S'_f(t) + S''_f(t),
\]
\[
S'_f(t) = \sum_{k \in \mathbb{Z}^n} f(k) \eta \left( \frac{|k|_Q}{t} \right), \tag{90}
\]
\[
S''_f(t) = \sum_{k \in \mathbb{Z}^n} f(k) (\chi_{[0,1]} - \eta) \left( \frac{|k|_Q}{t} \right). \tag{91}
\]

We can assume without loss of generality that \( f \geq 0 \). We begin with the first sum in \( S'_f(t) \):

**Lemma 5.3.**
\[
S'_f(t) = t^n \int_{\mathbb{R}^n} f(X) \eta(|X|_Q) dX + O(\log t).
\]

**Proof.** We write the sum as
\[
S'_f(t) = \sum_{j=0}^{\infty} \sum_{k \in \mathbb{Z}^n} f(k) \rho \left( \frac{2^j |k|_Q}{t} \right).
\]

We further break up the dyadic sum into \( \sum_{j=0}^{J(t)} + \sum_{j=J(t)+1}^{\infty} \) with \( J(t) \) to be determined later. We first consider
\[
S'_1 := \sum_{j=0}^{J(t)} \sum_{k \in \mathbb{Z}^n} f(k) \rho \left( \frac{2^j |k|_Q}{t} \right).
\]
The function $f(X)\rho(2^j|X|_Q) \in C_0^\infty(\mathbb{R}^n)$ when $f$ is homogeneous of degree 0 and smooth on $\partial Q$. Hence we may apply the Poisson summation formula to the $k$ sum to obtain

$$S'_1 = \sum_{j=0}^{J(t)} \sum_{N \in \mathbb{Z}^n} \int_{\mathbb{R}^n} e^{-i\langle X, N \rangle} f(X) \rho \left( \frac{2^{j}|X|_Q}{t} \right) dX.$$

The terms with $N = 0$ sum up to

$$t^n \int_{\mathbb{R}^n} f(X) \left[ \sum_{j=0}^{J(t)} \rho(2^j|X|_Q) \right] dX = t^n \int_{\mathbb{R}^n} f(X) \left[ \eta(|X|_Q) - \eta(2^{J(t)+1}|X|_Q) \right] dX = t^n \int_{\mathbb{R}^n} f(X) \eta(|X|_Q) dX + O(t^n 2^{-nJ(t)}),$$

where the last estimate is a consequence of the fact that $\eta(2^{J(t)+1}|X|_Q)$ is supported on $2^{-J(t)} Q$.

To estimate the remaining terms in the sum $S'_1$, we make the change of variables $Y = 2^j X/t$ in the integral to obtain

$$2^{-nj} t^n \int_{\mathbb{R}^n} f(Y) \rho(|Y|) e^{-i2^{-j}t\langle Y, N \rangle} dY.$$

Since the integrand is smooth, this term has the upper bound

$$c 2^{-nj} t^n (1 + 2^{-j}|N|t)^{-K}, \quad \forall K > 0.$$

(Again, we let $c$ denote a constant; $c$ depends on $f$ and $K$, but is independent of $j$, $t$, $N$.) The sum over $N \neq 0$ is then bounded by

$$c t^n \sum_{j \leq J(t)} 2^{-nj} \sum_{N \neq 0} (1 + 2^{-j}|N|t)^{-K} \sim t^n \sum_{j \leq J(t)} 2^{-nj} \int_0^{\infty} (1 + 2^{-j}rt)^{-K} r^{n-1} dr = \sum_{j \leq J(t)} \int_0^{\infty} (1 + s)^{-K}s^{n-1} ds = c J(t).$$

Therefore

$$S'_1 = t^n \int_{\mathbb{R}^n} f(X) \eta(|X|_Q) dX + O(t^n 2^{-nJ(t)}) + O(J(t)).$$

Recall that $S'_f(t) = S'_1 + S'_2$, where

$$S'_2 = \sum_{j=J(t)+1}^{\infty} \sum_{k \in \mathbb{Z}^n} f \left( \frac{k}{|k|_Q} \right) \rho \left( \frac{2^{j}|k|_Q}{t} \right).$$

Since

$$\sum_{j=J(t)+1}^{\infty} \rho \left( \frac{2^{j}|k|_Q}{t} \right) = \eta \left( \frac{2^{j(t)}|k|_Q}{t} \right) \leq \chi t 2^{-j(t)} Q,$$
the remainder $S'_2$ is bounded by the total number of lattice points in the shell $|k|_Q \leq 2^{-J(t)t}$, hence is of order $t^n 2^{-nJ(t)}$. It follows that
\[
S'_f(t) = t^n \int_{\mathbb{R}^n} f(X) \eta(|X|_Q) dX + O(t^n 2^{-nJ(t)}) + O(J(t)).
\] (92)

To balance the terms, we choose $J(t) = \log_2 t$, and then the last two terms of (92) have the form
\[
O(t^n t^{-n}) + O(\log t) = O(\log t).
\]

\[\square\]

5.1.2. Stationary phase. Theorem 5.1 is an immediate consequence of Lemma 5.3 and the following asssymptotics of the second sum $S''_f(t)$ from (91):

**Lemma 5.4.**
\[
S''_f(t) = t^n \int_{\mathbb{R}^n} f(X)(\chi_{[0,1]} - \eta)(|X|_Q) dX + O\left(t^n \frac{2n}{n+1}\right).
\]

**Proof.** Let
\[
g(X) = f(X)(\chi_{[0,1]} - \eta)(|X|_Q)
\]
and mollify $g$ by a radial approximate identity $\varphi_\varepsilon$ to obtain a smooth approximation $g_\varepsilon = g \ast \varphi_\varepsilon$. We claim that
\[
S''_f(t) = \sum_{k \in \mathbb{Z}^n} g\left(\frac{k}{t}\right) = \sum_{k \in \mathbb{Z}^n} g_\varepsilon\left(\frac{k}{t}\right) + O(\varepsilon t^n).
\] (93)

To see this, we break the sum into two parts. The first part is over the lattice points $k$ with $k/t$ in an $\varepsilon$ tube $T_\varepsilon$ about $\{|X|_Q = 1\}$. The number of such $k$ is $O(\varepsilon t^n)$, so this part contributes the stated error. For the remaining sum, the error is
\[
\left| \sum_{k \in \mathbb{Z}^n \setminus T_\varepsilon} \left[ g\left(\frac{k}{t}\right) - g_\varepsilon\left(\frac{k}{t}\right) \right] \right| \leq \sum_{k/t \in \text{Supp} g \setminus T_\varepsilon} \varepsilon \sup_{|X|_Q > 1} |dg(X)| = O(\varepsilon t^n),
\]
which verifies (93).

The Poisson summation formula then gives
\[
\sum_{k \in \mathbb{Z}^n} g_\varepsilon(k/t) = t^n \sum_{N \in \mathbb{Z}^n} \hat{g}_\varepsilon(2\pi t N) = t^n \sum_{N \in \mathbb{Z}^n} \hat{g}(2\pi t N) \hat{\varphi}(2\pi t N).
\]

The term $N = 0$ yields
\[
t^n \int_{\mathbb{R}^n} g_\varepsilon(X) dX = t^n \int_{\mathbb{R}^n} f(X)(\chi_{[0,1]} - \eta)(|X|_Q) dX + O(\varepsilon t^n),
\]
where the last inequality is by breaking up the integral into two parts as above.

As for the remainder terms $N \neq 0$, we now show that
\[
\hat{g}(2\pi t N) = O\left((1 + |tN|)^{-\frac{n+1}{2}}\right).
\] (94)
To verify (94), we write
\[ g = -f \rho h = -(f \rho)(h \eta_2), \]
with \( \eta_2(X) = \eta \left( \frac{1}{2} |X|_Q \right), \) \( h = \theta \circ \lambda, \)
\( \lambda(X) = |X|_Q - 1, \) \( \theta(t) = \text{Heaviside function} = \begin{cases} 0, & \text{if } t < 0, \\ 1, & \text{if } t \geq 0. \end{cases} \)

Since \( \widehat{g} = -\widehat{f \rho} \ast \widehat{h \eta_2} \) and \( \widehat{f \rho} \) is rapidly decaying, it suffices to show that \( \widehat{h \eta_2} \) satisfies (94). (Here, we use the elementary estimate \( \| \alpha \ast \beta \|_K \leq c \| \alpha \|_{(K+1)} \| \beta \|_K, \) where \( \| \alpha \|_K = \sup_{x \in \mathbb{R}^n} (1 + |x|)^K |\alpha(x)|. \) Taking partial derivatives,
\[
D_j(h \eta_2) = D_j \eta_2 + (\delta_0 \circ \lambda) D_j \lambda.
\]
Since the latter term is given by integration over \( \partial Q, \) which is strictly convex, the standard stationary phase method (see Hörmander [Ho]) immediately gives
\[
\widehat{\delta_0 \circ \lambda}(x) = O \left( x^{-\frac{n-1}{2}} \right),
\]
which implies (94).

Hence the remainder is bounded above by
\[
c t^n \sum_{N \neq 0} (1 + |tN|)^{-(n+1)/2} (1 + |\varepsilon t N|)^{-K}. \tag{95}
\]

The sum (95) can be replaced by the integral
\[
c t^n \int_{\mathbb{R}^n} (1 + |tN|)^{-(n+1)/2} (1 + |\varepsilon t N|)^{-K} dN = c t^n \int_0^\infty (1+tr)^{-(n+1)/2} (1+\varepsilon tr)^{-K} r^{n-1} dr
\]
\[
= c \varepsilon^{1-n/2} \int_0^\infty (\varepsilon+s)^{-(n+1)/2} (1+s)^{-K} s^{n-1} ds
\]
\[
\leq c \varepsilon^{1-n/2} \int_0^\infty (1+s)^{-K} s^{n-3} ds = c \varepsilon^{1-n/2}. \]

Hence
\[
S'_f(t) = t^n \int_{\mathbb{R}^n} f(X)(X(0,1] - \eta) |X|_Q dX + O(\varepsilon t^n) + O(\varepsilon^{-(n-1)/2}).
\]

To optimize, we choose \( \varepsilon \) so that \( \varepsilon t^n = \varepsilon^{-(n-1)/2}, \) i.e. \( \varepsilon = t^{2n/(n+1)}, \) which gives the result. (To be precise, it is the sum of the terms in (95) with \( |N| \geq \sqrt{n} \) that is bounded by the above integral. But there are only finitely many terms with \( |N| < \sqrt{n}, \) and each of these terms is \( < c t^{n-\frac{n+1}{2}}, \) which is better than \( O \left( t^{n-\frac{2n}{n+1}} \right) \) when \( n \geq 2. \) \( \square \)
5.1.3. Van der Corput for homogeneous weights $f$. Proof of Corollary 5.2. This time, we have

$$S_f(t) = \sum_{k \in \mathbb{Z}^n \cap Q < [0]} |k| \cdot f \left( \frac{k}{|k|} \right).$$

The set of norms of lattice points $\{t_j \in \mathbb{R}^+ : \exists k \in \mathbb{Z}^n \ni |k|_Q = t_j\}$ is a countable set without accumulation point. We order the $t_j$ so that $t_j \leq t_{j+1}$. We then define the monotone increasing step function on $\mathbb{R}$,

$$\mu(T) = \sum_{j : t_j \leq T} \left\{ \sum_{k : |k|_Q = t_j} f \left( \frac{k}{|k|} \right) \right\}.$$

It is clear that

$$\mu(T) = S_{f_0}(T), \quad f_0(x) = \frac{f(x)}{|x|_Q}.$$

Hence, by Theorem 5.1,

$$S_{f_0}(t) = t^n \int_Q f_0 \, dX + O \left( t^{n - \frac{2n}{n+1}} \right), \quad t \to \infty. \quad (96)$$

We further have

$$S_f(T) = \int_0^T t^\alpha \, d\mu(t). \quad (97)$$

Indeed,

$$d\mu(t) = \sum_j \left\{ \sum_{k : |k|_Q = t_j} f \left( \frac{k}{|k|} \right) \right\} \delta(t_j),$$

and

$$\int_0^T t^\alpha \, d\mu(t) = \sum_{j : t_j \leq T} \left\{ \sum_{k : |k|_Q = t_j} f \left( \frac{k}{|k|} \right) \right\} t_j^\alpha = S_f(T).$$

Integrating (97) by parts and applying (96), we get

$$S_f(T) = T^\alpha \mu(T) - \alpha \int_0^T t^{\alpha-1} \mu(t) \, dt$$

$$= T^\alpha \left[ T^n \int_Q f_0 \, dX + O \left( T^{n - \frac{2n}{n+1}} \right) \right] - \alpha \int_0^T t^{\alpha-1} \left[ \int_Q f_0 \, dX + O \left( T^{n - \frac{2n}{n+1}} \right) \right] \, dt$$

$$= T^{n+\alpha} \left[ \int_Q f_0 \, dX \right] \frac{n}{\alpha + n} + O \left( T^{n - \frac{2n}{n+1} + \alpha} \right)$$

$$= T^{n+\alpha} \int_Q f \, dX + O \left( T^{n - \frac{2n}{n+1} + \alpha} \right).$$
5.2. Van der Corput for critical points. We prove Theorem 1.4 by following the arguments of the proofs of Theorem 5.1 and Corollary 5.2 with hardly any changes. We first assume that $\psi$ is homogeneous of order 0 in $S$. We let $K_\psi = \rho (\text{Supp } \psi) \subset C$, a compact set.

To begin, we recall that if $W$ has critical points, then $W$ is in the ‘light cone’ $Q[W] > 0$. For $W$ in the light cone, we write

$$|W|_Q = Q[W]^{\frac{1}{2}}, \quad \text{for } Q[W] > 0.$$  

The main difference between this case and our previous one, is that now the set $Q$ given by (89), in addition to not being convex, is not compact. However, since only with those $W$ with critical points in the support of $\psi$ contribute to the sum, we consider

$$Q_\psi := Q \cap S_\psi, \quad S_\psi = \left( \bigcup_{\tau \in K_\psi} S_\tau \right),$$

which is a compact subset of $S$.

We let $f(W) = \langle C_W, \psi \rangle$, which is a smooth function supported in $S_\psi$. Then

$$N_\psi(L) = S_f(L) = \sum_{k \in \mathbb{Z}^n \cap \sqrt{L} Q \setminus \{0\}} f(k),$$

as before. Now we follow the previous proof, with $t = \sqrt{L}$. Our first modification is to verify (93), we instead let $T_\varepsilon$ be the epsilon tube over $S_\psi \cap \partial Q$. Finally, the estimate $(\delta_0 \circ \lambda) (t) = O \left( t^{-\frac{n-1}{2}} \right)$, which was based on the convexity of $Q$ in our previous argument, holds in this case, since the phase $\psi(Y) = L \langle Y, N \rangle$ has (two) non-degenerate critical points whenever $N$ is in the light cone. Thus we have

$$N_\psi(L) = L^{b_3} \left[ \int_{\{Q[W] \leq 1\}} \langle C_W, \psi \rangle \, dW + O \left( L^{-\frac{2b_3}{2b_3+1}} \right) \right].$$

The case $\alpha = 0$ now follows from Lemma 4.3, and the general case then follows exactly as in the proof of Corollary 5.2. \qed

6. Special Geometry and Density of Critical Points

The aim of this section is to compute the critical point density $k^{\text{crit}}(Z)$ and verify Corollaries 1.5–1.6. At the same time, we compute the index density and prove Theorem 1.8. As in [DSZ1], we do this by pushing forward the integrand of (25) under the Hessian map. The Hessian map turns out to be an isomorphism, hence the discussion is more elementary than in [DSZ1]. To make the change of variables, we first evaluate the image of the Hessian using the special geometry of Calabi-Yau moduli spaces and then check how the Hessian map distorts inner products. Our discussion gives an alternate approach to the formulas in the article [DD1], and connects the special critical point density formula in this article with the general ones in [DSZ1, DSZ2].
6.1. The range of the Hessian map. We now study the complex Hessian map:

\[ H^c(Z) : W \to \begin{pmatrix} H' & -x \otimes (Z) \\ -\bar{x} \otimes (Z) & \bar{H}' \end{pmatrix}. \] (98)

To describe \( H^c(Z) \) in local coordinates, we fix a point \( Z_0 = (z_0, \tau_0) \) and choose normal coordinates \( \{z^1, \ldots, z^{h^{2,1}}\} \) at \( z_0 \in \mathcal{M} \). We let \( \Omega \) be a local normal frame for \( H^3,0 \to \mathcal{M} \) at \( z_0 \), and we let \( \omega = dx + \tau dy \). Recall that \( \omega \) is not a normal frame, since \( |\omega| = (\text{Im} \tau)^{1/2} \). We let \( \bar{\mathcal{E}}_\lambda = (\text{Im} \tau_0)^{1/2} \Omega^* \otimes \omega^* \), so that \( |\bar{\mathcal{E}}_\lambda(Z_0)| = 1 \).

As in §3.2, the matrix \( (H_{jk}) \) of the holomorphic Hessian is given by

\[ H'(Z_0) = \sum_{j,q} H'_{jq} dz^j \otimes dz^q \otimes \bar{\mathcal{E}}_\lambda|_{Z_0}, \quad 0 \leq j, q \leq h^{2,1}, \] (99)

where

\[ dz^0|_{Z_0} = \frac{1}{\text{Im} \tau_0} d\tau|_{Z_0} \]

is the unit holomorphic cotangent vector (with respect to the Weil-Petersson, or hyperbolic, metric on \( \mathcal{E} \)) at \( \tau_0 \).

We wish to express formulas (59)–(60) for the complex Hessian in terms of these coordinates and frames. We write

\[ (\nabla_j f) \otimes e_\lambda = \nabla_{\partial_j \partial \lambda}(fe_\lambda), \quad 1 \leq j \leq h^{2,1}, \quad (\nabla_0 f) \otimes e_\lambda = (\text{Im} \tau_0)\nabla_{\partial_0 \partial \lambda}(fe_\lambda). \]

\((\nabla_0 \) is the normalized covariant \( \tau \)-derivative given by (87).) The complex Hessian matrix is given by:

\[ H^c(Z_0) = \begin{pmatrix} H'(Z_0) & f(Z_0) I \\ f(Z_0) I & H'(Z_0) \end{pmatrix}, \quad H' = \left( \nabla_j \nabla_q f \right)_{0 \leq j, q \leq h^{2,1}}. \] (100)

Identifying the off-diagonal components with \( f(Z_0) \in \mathbb{C} \), we view the image space as a subspace of \( \text{Sym}\langle h^{2,1} + 1, \mathbb{C} \rangle \oplus \mathbb{C} \), so we can write the Hessian map in the form

\[ H_{Z_0} : S_Z \to \text{Sym}\langle h^{2,1} + 1, \mathbb{C} \rangle \oplus \mathbb{C}, \quad W \mapsto (H'(Z_0), f(Z_0)). \]

Lemma 6.1. The range of the Hessian map \( H_{Z_0} : S_Z \to \text{Sym}\langle h^{2,1} + 1, \mathbb{C} \rangle \oplus \mathbb{C} \) is of the form \( \mathcal{H}_{Z_0} \oplus \mathbb{C} \), where \( \mathcal{H}_{Z_0} \) is a real subspace of \( \text{Sym}\langle h^{2,1} + 1, \mathbb{C} \rangle \) of real dimension \( 2h^{2,1} \) spanned over \( \mathbb{R} \) by the matrices

\[ \xi^k = \begin{pmatrix} 0 & e_k \\ e_k^* & \mathcal{F}^k(z) \end{pmatrix}, \quad \xi^{h^{2,1}+k} = \begin{pmatrix} 0 & \sqrt{-1} e_k \\ \sqrt{-1} e_k^* & -\sqrt{-1} \mathcal{F}^k(z) \end{pmatrix}, \quad 1 \leq k \leq h^{2,1}, \]

given by (27), where \( e_k \) is the \( k^{th} \) standard basis element of \( \mathbb{C}^{h^{2,1}} \) and \( \mathcal{F}^k(z) \in \text{Sym}\langle h^{2,1}, \mathbb{C} \rangle \) is the matrix \( \mathcal{F}_{jk}^k(z) \) of (46).
In other words, $\mathcal{H}_{Z_0}$ is the set of matrices of the form
\[
\begin{pmatrix}
0 & (\tilde{v}_1, \ldots, \tilde{v}_{h^{2,1}}) \\
(\tilde{v}_1, \ldots, \tilde{v}_{h^{2,1}})' & \sum_{k=1}^{h^{2,1}} F^k(z) v_k
\end{pmatrix},
\]
where $(v_1, \ldots, v_{h^{2,1}}) \in \mathbb{C}^{h^{2,1}}$. (101)

We emphasize that $\mathcal{H}_Z \subset \text{Sym}(\mathbb{H}^{2,1} + 1, \mathbb{C})$ is only a real and not a complex subspace. We also note that $\dim_{\mathbb{R}} \mathcal{H}_Z = 2h^{2,1}$ and hence $\dim_{\mathbb{R}} (\mathcal{H}_Z \oplus \mathbb{C}) = b_3 = \dim_{\mathbb{R}} S_Z$; i.e., $\mathcal{H}_Z$ is an isomorphism.

**Proof of Lemma 6.1.** We shall use the notation $1 \leq j, k, l \leq h^{2,1}, 0 \leq \alpha, \beta, \gamma \leq h^{2,1}$. By (74), we have the (real-linear) isomorphism
\\[
\tilde{\mathcal{W}}_{Z_0} = \mathcal{W} \circ I^{-1}_r : H_{z_0}^{2,1} \oplus H_{z_0}^{0,3} \approx S_{Z_0}.
\]
Recall from proposition 2.2 that $H_{z_0}^{2,1} \oplus H_{z_0}^{0,3}$ has a complex $Q_{Z_0}$-orthonormal basis $\{\chi_{\alpha}\}$ of the form
\[
\chi_j = D_j \Omega_{Z_0}, \quad 1 \leq j \leq h^{2,1}, \quad \chi_0 = \overline{\Omega}_{Z_0}.
\]
By (76), a real $Q_{Z_0}$-orthonormal basis of $S_{Z_0}$ is
\[U_{\alpha} := (\text{Im } \tau)^{1/2} \tilde{\mathcal{W}}_{Z_0}(\chi_{\alpha}), \quad V_{\alpha} := (\text{Im } \tau)^{1/2} \tilde{\mathcal{W}}_{Z_0}(\sqrt{-1} \chi_{\alpha}) .
\]
In the above local frame $\tilde{e}_\mathcal{L}$; we write
\[U_{\alpha} = f_{\alpha} \tilde{e}_\mathcal{L}, \quad V_{\alpha} = g_{\alpha} \tilde{e}_\mathcal{L} ;
\]
equivalently, in the local $e_\mathcal{L}$ of Section 3.3
\[
\tilde{\mathcal{W}}_{Z_0}(\chi_{\alpha}) = f_{\alpha} e_\mathcal{L}, \quad \tilde{\mathcal{W}}_{Z_0}(\sqrt{-1} \chi_{\alpha}) = g_{\alpha} e_\mathcal{L} .
\]
We must compute the matrices
\[H'_{Z_0}(f_{\alpha} \tilde{e}_\mathcal{L}) = (\nabla_\beta \nabla_\gamma f_{\alpha})|_{Z_0}, \quad H'_{Z_0}(g_{\alpha} \tilde{e}_\mathcal{L}) = (\nabla_\beta \nabla_\gamma g_{\alpha})|_{Z_0},
\]
where $H'_{Z_0} : S_{Z_0} \rightarrow \text{Sym}(\mathbb{H}^{2,1} + 1, \mathbb{C})$ is the holomorphic Hessian map.
We shall show that:
\[
\begin{align*}
(i) \quad & \nabla^2_0 f_G(Z_0) = 0, \quad \forall G \in H^3_{Z_0}(X, \mathbb{C}) \quad (\text{where } W_G = f_G e_\mathcal{L})
\quad \text{and thus } \nabla^2_0 f_{\alpha}(Z_0) = \nabla^2_0 g_{\alpha}(Z_0) = 0, \\
(ii) \quad & \nabla_j \nabla_0 f_0(Z_0) = \nabla_j \nabla_0 g_0(Z_0) = 0, \\
(iii) \quad & \nabla_k \nabla_j f_0(Z_0) = \nabla_k \nabla_j g_0(Z_0) = 0, \\
(iv) \quad & \nabla_k \nabla_0 f_j(Z_0) = -\sqrt{-1} \delta_{jk}, \quad \nabla_k \nabla_0 g_j(Z_0) = -\delta_{jk}, \\
(v) \quad & \nabla_k \nabla_l f_j(Z_0) = F^j_{kl}, \quad \nabla_k \nabla_l g_j(Z_0) = \sqrt{-1} F^j_{kl},
\end{align*}
\]
First,
\[
\nabla_0 f_G(z, \tau) = \frac{|\text{Im } \tau_0|}{\text{Im } \tau} \int_X (F + \bar{\tau} H) \wedge \Omega_z.
\]
It follows that
\[ \nabla_0^2 f_G(z_0, \tau_0) = \frac{|\text{Im} \tau_0|^2}{\text{Im} \tau} \frac{\partial}{\partial \tau} \int_X (F + \bar{t} H) \wedge \Omega_z = 0 \]
by the critical point equation \( \nabla_0 f_G(z_0, \tau_0) = 0 \). This proves (i).

Next, differentiating (103) with \( f_G = f_\alpha \), we get
\[ \nabla_j \nabla_0 f_\alpha(Z_0) = \int \overline{\chi_\alpha} \wedge D_j \Omega Z_0 = \int \overline{\chi_\alpha} \wedge \chi_j = -i \delta_{j\alpha}, \]
and similarly,
\[ \nabla_j \nabla_0 g_\alpha(Z_0) = \int i \chi_\alpha \wedge \chi_j = -\delta_{j\alpha}. \]
This verifies (ii) and (iv).

Finally, we have by (46),
\[ \nabla_k \nabla_j f_\alpha = \int \chi_\alpha \wedge D_k D_j \Omega = -i \sum_l F^i_{kj} \int \chi_\alpha \wedge \overline{D_l \Omega}, \]
and hence
\[ \nabla_k \nabla_j f_\alpha(Z_0) = -i \sum_l F^i_{kj} \int \chi_\alpha \wedge \overline{\chi_l} = -i \sum_l F^i_{kj} \delta_{l\alpha} = \begin{cases} -i F^i_{kj} & \text{for } \alpha \geq 1 \\ 0 & \text{for } \alpha = 0 \end{cases} . \]
We also have \( \nabla_k \nabla_j g_\alpha(Z_0) = i \nabla_k \nabla_j f_\alpha(Z_0) \), verifying (iii) and (v).

Thus, the holomorphic Hessian \( H'(Z_0) \) maps the orthonormal fluxes
\[ i U_1, \ldots, i U_{h^2,1}, -i V_1, \ldots, -i V_{h^2,1} \]
given by (27). Furthermore,
\[ f_0(Z_0) = 1, \ H'(U_0) = 0, \ g_0(Z_0) = i, \ H'(V_0) = 0, \]
while
\[ f_j(Z_0) = g_j(Z_0) = 0. \]
Thus \( H^c(Z_0) \) maps the orthonormal fluxes (104) to the elements \( \xi^a \oplus 0 \in \text{Sym}(h^{2,1} + 1, \mathbb{C}) \oplus \mathbb{C} \), and maps \( U_0 \) to \( 0 \oplus 1 \) and \( V_0 \) to \( 0 \oplus i \). □
6.2. Distortion of inner product under the Hessian map. We recall that the space $\text{Sym}(h^{2,1} + 1, \mathbb{C})$ of complex symmetric matrices, regarded as a real vector space, has the inner product

$$(A, B)_{\mathbb{R}} = \text{Re} \langle A, B \rangle_{HS} = \text{Re} (\text{Trace } AB^*) .$$

(105)

Recalling that $S_Z = \tilde{\mathcal{N}}_Z(H^2_2, 1 + 1, H^0_2)$, we consider its codimension 1 subspace $S'_Z = \tilde{\mathcal{N}}_Z(H^2_2)$. By the proof of Lemma 6.1, the holomorphic Hessian map

$$H_Z : S'_Z \rightarrow \mathcal{H}_Z$$

(106)

is bijective, but as a map between inner product spaces, it is not an isometry. The distortion is given by the positive definite operator $\Lambda_Z$. We write

$$\Lambda_Z \xi^a = \sum_{b=1}^{2h^{2,1}} \Lambda_{ab} \xi^b ,$$

so that

$$(\xi^a, \xi^b)_{\mathbb{R}} = (\Lambda_Z^{-1} \Lambda_Z \xi^a, \xi^b)_{\mathbb{R}} = \sum_c \Lambda_{ac} (\Lambda_Z^{-1} \xi^c, \xi^b)_{\mathbb{R}} = \sum_c \Lambda_{ac} \delta_{cb} = \Lambda_{ab} .$$

Tracing through the definitions, we obtain that $(\Lambda_{ab})$ is the matrix

$$
\left(\begin{array}{cc}
\Lambda' & \Lambda'' \\
\Lambda'' & \Lambda'
\end{array}\right), \quad \Lambda'_{jk} = 2\delta_{jk} + \text{Re } \text{Tr } \mathcal{F}^j \mathcal{F}^{k*}, \quad \Lambda''_{jk} = \text{Im } \text{Tr } \mathcal{F}^j \mathcal{F}^{k*}
$$

(107)

of Hilbert-Schmidt inner products of the matrices in Lemma 6.1. Hence,

$$\Lambda'_{jk} + \sqrt{-1} \Lambda''_{jk} = 2\delta_{jk} + \text{Tr } \mathcal{F}^j \mathcal{F}^{k*} .$$

(108)

To tie this discussion together with that in [AD] and [DSZ2, §2.1], we note that we can consider $\mathcal{H}_Z$ as a complex vector space by redefining complex multiplication in $\mathcal{H}_Z$:

$$c \odot \begin{pmatrix} 0 & u \\ u^t & A \end{pmatrix} = \begin{pmatrix} 0 & \bar{c}u^t \\ \bar{c}u^t & cA \end{pmatrix} .$$

We then define a Hermitian inner product on $\mathcal{H}_Z$:

$$\left(\begin{pmatrix} 0 & u \\ u^t & A \end{pmatrix}, \begin{pmatrix} 0 & v \\ v^t & B \end{pmatrix}\right) = 2\bar{u} \cdot v + \text{Tr}(AB^*) .$$

The matrices $\{\xi_j, \ j = 1 \ldots, h^{2,1}\}$ form a complex basis of $\mathcal{H}_Z$. It follows from (29) that

$$\Lambda_Z = \sum_{j=1}^{h^{2,1}} \xi_j \otimes \xi_j^* .$$

(109)
where the $\xi^j$ are $(h^{2,1} + 1) \times (h^{2,1} + 1)$ matrices. Each term $\xi^j \otimes \xi^{j*}$ in $\Lambda Z$ may be expressed in matrix form as $\begin{pmatrix} \xi^j_{ab} \bar{\xi}^{j*}_{cd} \end{pmatrix}$; i.e.,

$$(\Lambda Z H)_{kl} = \sum_{p,q} [\Lambda Z]^{pq}_{kl} H_{pq}, \quad [\Lambda Z]^{pq}_{kl} = \sum_{j=1}^{h^{2,1}} \xi^j_{kl} \bar{\xi}^{j*}_{pq}, \quad 0 \leq k, l, p, q \leq h^{2,1}. \quad (110)$$

As in [Dsz2, 2.1] $\Lambda Z$ may be expressed in terms of the Szegö kernel $\Pi^Z$, i.e. the kernel of the orthogonal projection onto $S^Z$. Let $F^Z_\xi$ denote the local representative of $\Pi^Z$ in the frame $\tilde{e}_L \otimes \tilde{e}_L^*$. Using (102) and (109), we have

$$[\Lambda Z]^{pq}_{kl} = \nabla_{\xi_k} \nabla_{\xi_l} \nabla_{\bar{\eta}_p} \nabla_{\bar{\eta}_q} F^Z_\xi (\xi, \eta)|_{\xi=\eta=Z}. \quad (111)$$

In addition, $\Lambda Z$ determines an operator $\tilde{\Lambda} Z$ on the space $\mathcal{H}^c$ of complex matrices of the form

$$H^c := \begin{pmatrix} H & xI \\ xI & \bar{H} \end{pmatrix}, \quad H \in \text{Sym}(h^{2,1}, \mathbb{C}), \quad (112)$$

defined by

$$\tilde{\Lambda} Z \begin{pmatrix} H & xI \\ xI & \bar{H} \end{pmatrix} = \begin{pmatrix} \Lambda Z H & xI \\ xI & \overline{\Lambda Z \bar{H}} \end{pmatrix}. \quad (113)$$

We now relate the $(1, 1)$-form $\omega_\Lambda$ of (31) and the operator $\Lambda$ to the curvature of the Weil-Petersson metric on $\mathcal{C}$.

**Proposition 6.2.** We have:

i) $[\Lambda Z]^{jq}_{j'q'} = -G^{q\bar{p}} R^j_{jp} (\bar{q}) + \delta^q_j \delta^q_{q'} + \delta^j_q \delta^q_{j'}$, where $R$ is the curvature tensor of the Weil-Petersson metric on $\mathcal{C}$;

ii) $\omega_\Lambda = (m + 3) \omega_{WP} + \text{Ric}(\omega_{WP})$, where Ric is the Ricci curvature $(1, 1)$ form of the Weil-Petersson metric of $\mathcal{M}$, i.e.

$$\text{Ric}(\omega_{WP}) = \frac{i}{2} \sum_{ij} \text{Ric}_{ij} dz^i \wedge d\bar{z}^j, \quad \text{Ric}_{ij} := -G^k_{\bar{k}j} R^i_{j\bar{k}\bar{\epsilon}}.$$

Thus, $\omega_\Lambda$ is the Hodge metric [Lu, Wa2].

**Proof.** To prove (i), it suffices to combine (110) and (52), raising and lowering indices as appropriate. (In (110), a normal frame at $Z$ is assumed.)

For (ii) we note that the $(1, 1)$-form

$$\omega_\Lambda = \frac{i}{2} \sum_{ij} \left[ 2\delta_{ij} + \text{Tr} \mathcal{F}^i (Z) \mathcal{F}^{j*} (Z) \right] dz^i \wedge d\bar{z}^j. \quad (114)$$

On the other hand, by (47),

$$\text{Ric}_{ij} = -G^k_{\bar{k}j} \left[ G_{ij} G_{k\bar{k}} + G_{i\bar{\epsilon}} G_{k\bar{\epsilon}} - \frac{1}{\Omega \wedge \Omega} \sum_{p,q} G^{p\bar{q}} \mathcal{F}_{ikp} \mathcal{F}_{j\bar{q}} \right]$$

$$= -(m + 1) G_{ij} + Tr \mathcal{F}^i \mathcal{F}^{j*}. \quad (115)$$

\[\square\]
Remark. To facilitate comparison with [AD, DSZ1], we note that our notational conventions are the same as in [DSZ1]. In [AD], the Szegö kernel \( \Pi_1 \) is denoted \( G_\mathcal{Z} \). The formulas in [AD] (4.8) are the same as (110), resp. Proposition 6.2(1). Also \( F_{ab|\bar{c}d} = \Lambda_{ab}^p G_{pc} G_{q\bar{d}} \). The coefficients \( F_{ab|\bar{c}d} \) in [AD] correspond to the off-diagonal blocks of \( \tilde{\Lambda} \).

6.3. Proof of Theorem 1.8. All but one of the ingredients of the proof are precisely the same as in Theorem 1.4. We first define the analogue of (25) and (32) for the signed sum:

\[
\text{Ind}(Z) := \int_{S_Z} \det H^c W(Z) \chi_{QZ} dW = \frac{1}{b_3! \sqrt{\det \Lambda_Z}} \int_{H^c \otimes \mathbb{C}} \det \left( H^* H - |x|^2 I \right) e^{-\left( \Lambda_Z^{-1} H^* H - |x|^2 \right)} H dx .
\]

(116)

By Lemma 4.2 and the proof of Lemma 4.3, we conclude that

\[
\text{Ind}_{\chi_K}(L) = L^{b_3} \left[ \int_K \text{Ind}(Z) d\text{Vol}_{WP} + O \left( L^{-1/2} \right) \right] .
\]

(117)

To complete the proof of Theorem 1.8, we evaluate the integral in (116):

Lemma 6.3. We have

\[
b_3! \text{Ind}(Z) d\text{Vol}_{WP} = \frac{\pi^{2m}}{2^m} c_m (T^{* (1,0)} (\mathcal{C}) \otimes \mathcal{L}, \omega_{WP} \otimes h^*_W) \left( \frac{\pi}{2} \right)^m \det (-R - \omega \otimes I) .
\]

Proof. This follows by a supersymmetric formula for the determinant, used in this context in [AD] and also in [BSZ2]. We briefly review the fermionic formalism referring to [BGV, BSZ2] for further details in a similar setting. Let \( M = \left( M^j_j \right) \) be an \( n \times n \) complex matrix. Then,

\[
\det M = \int_{B^{2n}} e^{-\langle M \eta, \bar{\eta} \rangle} d\eta , \quad \langle M \eta, \bar{\eta} \rangle = \sum_{j,j'} \eta_j M^j_j \bar{\eta}_{j'},
\]

(118)

where \( \eta_j, \bar{\eta}_j \) (1 ≤ \( j \) ≤ \( n \)) are anti-commuting (or “fermionic”) variables. The integral \( \int_{B} = \int_{B^{2n}} \) is the Berezin integral, a notation for the linear functional \( \int_{B} : \Lambda^* \mathbb{C}^{2n} \to \mathbb{C} \) defined by

\[
\int_{B} |_{\Lambda^t \mathbb{C}^{2n}} = 0 \quad \text{for} \quad t < 2n , \quad \int_{B} (\prod_j \bar{\eta}_j \eta_j) = 1 .
\]

We now apply this formalism to \( \det (H^* H - |x|^2 I) = \det H^c \), where \( H^c \) is defined as in (112) and refer to the discussion in §6.2. The matrix \( H^c \) is of rank \( b_3 \), and we write

\[
\det H^c = \int_{B^{2b_3}} e^{-\langle H^c(\eta, \bar{\eta}), (\theta, \bar{\theta}) \rangle} d\eta d\theta ,
\]

(119)
where \( \eta = (\eta_1, \ldots, \eta_{b_3/2}) \), \( \theta = (\theta_1, \ldots, \theta_{b_3/2}) \), and
\[
\langle H^c(\eta, \bar{\eta}), (\theta, \bar{\theta}) \rangle = \sum (H_{jk}\eta_j\theta_k + x\delta_{jk}\eta_j\bar{\theta}_k + x\delta_{jk}\bar{\eta}_j\theta_k + \bar{H}_{jk}\bar{\eta}_j\bar{\theta}_k).
\]

The quadratic form \( (\Lambda_Z^{-1}H, H)_{\mathbb{R}} + |x|^2 \) in the exponent of the Gaussian integral may be expressed in the form \( \frac{1}{2}(\tilde{\Lambda}_Z^{-1}H^c, H^c) \), where \( \tilde{\Lambda}_Z \) is the restriction of the operator defined in (113) to \( H^c_Z \). Indeed, both quadratic forms are equivalent to \( Q_Z(W, W) \) under a linear change of variables \( W \rightarrow H_Z(W) \) in the case of \( \Lambda_Z \) and \( W \rightarrow H^c(W) \) in the case of \( \tilde{\Lambda}_Z \).

Then
\[
b_3! \text{Ind}(Z) = \frac{1}{\sqrt{\det \Lambda_Z}} \int_{H^c_Z} \int_{B^{2b_3}} e^{-\langle H^c(\eta, \bar{\eta}), (\theta, \bar{\theta}) \rangle} - \langle \Lambda_Z^{-1}H^c, H^c \rangle dH^c d\eta d\theta. \tag{120}
\]

We let
\[
\Omega = (\eta, \bar{\eta}) \otimes (\theta, \bar{\theta})' = \begin{pmatrix} (\eta_j\theta_k) & (\eta_j\bar{\theta}_k) \\ (\bar{\eta}_j\bar{\theta}_k) & (\bar{\eta}_j\theta_k) \end{pmatrix},
\]
so that \( \langle H^c(\eta, \bar{\eta}), (\theta, \bar{\theta}) \rangle = \langle H^c, \Omega \rangle = \text{Tr} H^c \Omega' \). Then the \( dH^c \) integral in (120) becomes the Fourier transform of the Gaussian function \( e^{-\langle \Lambda^{-1}H^c, H^c \rangle} \) evaluated at \( i\Omega \). Recalling that the Fourier transform of \( e^{-\langle Ax, x \rangle/2} \) equals \( (2\pi)^{n/2}/(\det A)^{-1/2} e^{-\langle A^{-1}x, x \rangle/2} \), we have that the \( dH^c \) integral equals \( (\det \tilde{\Lambda})^{1/2} e^{-1/4(\tilde{\Lambda}\Omega, \Omega)} \). After cancelling (det \( \tilde{\Lambda} \))
\[
b_3! \text{Ind}(Z) = \pi^m \int_{B^{2b_3}} e^{-\frac{1}{4}(\tilde{\Lambda}\Omega, \Omega)} d\eta d\theta, \tag{121}
\]

where in normal coordinates, we have (by (113) and Proposition 6.2)
\[
(\tilde{\Lambda}_Z\Omega, \Omega)_{\mathbb{R}} = \text{Trace} \left[ \begin{pmatrix} \Lambda_Z \eta \otimes \theta & \eta \otimes \bar{\theta} \\ \bar{\eta} \otimes \theta & \tilde{\Lambda}_Z \bar{\eta} \otimes \bar{\theta} \end{pmatrix} \right]^* \\
= \sum_{jqjq'} (\Lambda_{jq} \eta_j \theta_q \bar{\eta}_j \bar{\theta}_q + \Lambda_{jq} \eta_j \bar{\theta}_q \bar{\eta}_j \bar{\theta}_q) + \sum_{jq} (\eta_j \bar{\theta}_q \eta_j \theta_q + \bar{\eta}_j \bar{\theta}_q \bar{\eta}_j \bar{\theta}_q) \\
= 2 \sum_{jqjq'} (\Lambda_{jq} \eta_j \theta_q \delta_{jq} - \delta_{jq} \delta_{jq'}) \eta_j \theta_q \bar{\eta}_j \bar{\theta}_q \\
= 2 \sum_{jqjq'} (\delta_{jq} \delta_{jq'} - \delta_{jq} \delta_{jq'}) \eta_j \theta_q \bar{\eta}_j \bar{\theta}_q.
\]

(Here we used the fact that \( \Lambda_{jq} = \Lambda_{jq}^{-1} \); see (110).) Thus
\[
b_3! \text{Ind}(Z) = \pi^m \int_{B^{2b_3}} e^{-\frac{1}{2}(\delta_{jq} \delta_{jq'} - \delta_{jq} \delta_{jq'}) \eta_j \bar{\eta}_j \bar{\theta}_q \theta_q} d\eta d\theta
\]
\[
= (\pi)^m \frac{\det (-R - \omega \otimes I)}{d\text{Vol}_W}.
\]

\( \Box \)

Remark. The index density computation in special geometry is closely related to the asymptotics in [DSZ2, §5] for critical point densities for powers of a positive line bundle \( L \) on a compact Kähler manifold \( M \). The expansions in §5.1 of [DSZ2] can be used to show that the (first few) terms in the asymptotic expansion of the index density equal those of the Chern form corresponding to \( c_m(T^{*1,0} \otimes L^N) \).
6.4. Examples. We describe in this section the critical point distribution for the cases where the dimension $h^{2,1}(X)$ of the moduli space is 0 and 1, i.e. when $\dim \mathcal{C}$ is 1 and 2, respectively.

6.4.1. $h^{2,1}(X) = 0$. The simplest example is the case where the Calabi-Yau manifold $X$ is rigid, i.e. $\mathcal{M} = \{ pt \}$. (See [AD, DD1] for further details and computer graphics of critical points in this case.) Then only the parameter $\tau \in \mathcal{H}$ varies. Let $G = F + i H$, and consider the flux superpotential $W_G$. Its critical point equation is

$$F + \tau H \in H^{0,3}$$

(since in this case $H^{2,1}(X, \mathbb{C}) = 0$). So we write

$$F = A \Omega + A \Omega \quad , \quad H = B \Omega + B \Omega \quad ,$$

$$A = a_1 + i a_2, \quad B = b_1 + i b_2 \in \mathbb{Z} + \sqrt{-1} \mathbb{Z} .$$

Then writing $W_G = W_{A,B}$, we have

$$\nabla W_{A,B} = 0 \iff F + \tau H \in H^{0,3} \iff A + \tau B = 0 \iff \tau = -\frac{A}{B} .$$

Each flux superpotential $W_{A,B} \in S$ (with $A, B \in \mathbb{C}$) has a unique critical point in $\mathcal{H}$, which may or may not lie in the fundamental domain $\mathcal{C}$. In the notation of (15),

$$\pi(S) = \{ W_{A,B} : -\frac{A}{B} \in \mathcal{C} \}$$

is a domain with boundary in $\mathbb{C}^2$. Each $SL(2, \mathbb{Z})$-orbit of fluxes (or superpotentials) contains a unique element whose critical point lies in $\mathcal{C}$, so $\pi(S)$ is a fundamental domain for the action of $\Gamma'$ on $S$.

Thus, counting critical points is equivalent to counting $SL(2, \mathbb{Z})$ orbits of superpotentials satisfying the tadpole constraint. The pair $(A, B)$ corresponds to the element

$$\begin{pmatrix} a_1 & b_1 \\ a_2 & b_2 \end{pmatrix} \in GL(2, \mathbb{Z})$$

and the Hodge-Riemann form quadratic form may be identified with the indefinite quadratic form

$$Q((A, B)) = a_1 b_2 - b_2 a_1$$

on $\mathbb{R}^4$. The modular group $SL(2, \mathbb{Z})$ acts by the standard diagonal action on $(A, B) \in \mathbb{R}^2 \times \mathbb{R}^2$ preserving $Q((A, B))$ or equivalently by left multiplication preserving det. Thus, the set of superpotentials satisfying the tadpole constraint is parametrized by:

$$\left\{ \begin{pmatrix} a_1 & b_1 \\ a_2 & b_2 \end{pmatrix} \in GL(2, \mathbb{Z}) : 0 < \det \begin{pmatrix} a_1 & b_1 \\ a_2 & b_2 \end{pmatrix} \leq L \right\} ,$$

and we want to count the number of $SL(2, \mathbb{Z})$-orbits in this set. Counting the number of $SL(2, \mathbb{Z})$ orbits in $D_L$ is equivalent to determining the average order of the classical divisor function $\sigma(m)$, see for instance Hardy-Wright [HW, Theorem 324]:

$$N_{\text{crit}}^L = \sum_{m=1}^{L} \sum_{k|m} k = \sum_{m=1}^{L} \sigma(m) \sim \frac{\pi^2}{12} L^2 + O(L \log L) . \quad (122)$$

As verified in [DD1] (and as follows very simply from Theorem 1.4), the critical points are uniformly distributed relative to the hyperbolic area form.
6.4.2. $h^{2,1}(X) = 1$. We now illustrate our notation and results with the case where the moduli space of complex structures on $X$ is one-dimensional over $\mathbb{C}$. (This case is also studied in [DD1] from a slightly different point of view.) In this case, there is a single Yukawa coupling $\mathcal{F}_{11}^1(z)$ defined by $D_2^2 \Omega_z = \mathcal{F}_{11}^1(z) \overline{D_z} \Omega_z$.

The space $S_{\zeta, \tau} \simeq H^{2,1} \oplus H^{0,3} \simeq \mathbb{C}^2$. The space is spanned as a real vector space by four superpotentials $U_0, U_1, V_0, V_1$ corresponding to $\{ \Omega_z, D_z \Omega_z, i \overline{\Omega}_z, i D_z \Omega_z \}$. By the proof of Lemma 6.1, the holomorphic Hessians of $U_0$ and $V_0$ at a critical point equal zero, so we only need to consider the holomorphic Hessian map on $U_1$ and $V_1$. The corresponding space of Hessians is the real 2-dimensional subspace $H_Z$ of $\text{Sym}(2, \mathbb{C})$ spanned by

$$\xi^1 = \begin{pmatrix} 0 & 1 \\ 1 & F(z) \end{pmatrix}, \quad \xi^2 = \sqrt{-1} \begin{pmatrix} 0 & 1 \\ 1 & -F(z) \end{pmatrix},$$

where we write $F = \mathcal{F}_{11}^1$. Hence, we may parameterize the space $H_Z$ of holomorphic Hessians by

$$w = y_1 + i y_2 \mapsto H(w) = \begin{pmatrix} 0 & w \\ w & F(z) \overline{w} \end{pmatrix}.$$

By (25), we have:

$$K_{\text{crit}}(Z) = \frac{1}{2} \int_{\mathbb{C} \oplus \mathbb{C}} | \det(H(w)^* H(w) - |x|^2 I) | \ e^{-|w|^2 + |x|^2} dw dx.$$

We note that

$$\det(H(w)^* H(w) - |x|^2 I) = |w|^4 + |x|^4 - (2 + |F(z)|^2) |x|^2 |w|^2.$$

Hence

$$K_{\text{crit}}(Z) = \frac{1}{2} \int_{\mathbb{C} \oplus \mathbb{C}} |w|^4 + |x|^4 - (2 + |F(z)|^2) |x|^2 |w|^2 \ e^{-|w|^2 + |x|^2} dw dx,$$

agreeing with (3.19) of [DD1]. There, the integral is evaluated as

$$K_{\text{crit}}(Z) = \frac{\pi^2}{2} \left( 2 - |F|^2 + \frac{2 |F|^3}{\sqrt{4 + |F|^2}} \right).$$

Remark. In this example, the discriminant variety is given by

$$\overline{D} = \{(Z, x W_0(Z) + w W_1(Z)) \in I : |w|^2 - |x|^2 = \pm |w x F(z)|^2 \},$$

where $W_\alpha = U_\alpha + i V_\alpha$. The matrix $\Lambda$ is given by

$$\Lambda = \begin{pmatrix} 2 + |F|^2 & 0 \\ 0 & 2 + |F|^2 \end{pmatrix}.$$
7. Problems and Heuristics on the String Theory Landscape

In this section, we continue the discussion begun in §1.6 on the bearing of our methods and results on the physicists’ picture of the string theory landscape. We briefly review some of the heuristic estimates in the physics discussions, and then discuss a number of mathematical pitfalls in the heuristics. In §7.2, we state some mathematical problems suggested by the heuristics and by rigorous vacuum statistics. In §7.3, we give our own (tentative) heuristic estimate of the dependence of the critical point density \( K_{\text{crit}}(Z) \) on the dimension \( b_3/2 \) of \( \mathcal{C} \).

7.1. Complexity of the string theory landscape. As mentioned in §1.6, the possible vacua in string/M theory are often represented as valleys in a complex string theory landscape, and the number of valleys is often estimated at \( 10^{500} \).

L. Susskind and others have argued that such a large number of possible vacua should essentially be a consequence of the large number of variables in the potential. A common and general argument to arrive at this number of vacua without specifying any particular string theory model is to reason that the potential energy is a function of roughly 1000 variables. A generic polynomial \( f \) of degree \( d \) on \( \mathbb{C}^m \) has \((d-1)^m\) critical points since critical points are solutions of the \( m \) equations \( \frac{\partial f}{\partial z_j}(w) = 0 \) of degree \( d-1 \). Thus, the number of critical points would seem to grow at an exponential rate in the number of variables. Such an exponential growth rate of critical points also appears in the physics of spin glasses, where the growth in the number of metastable states (local minima of the Hamiltonian) in terms of the number of variables is often used to measure the complexity of the energy landscape. In the special model of random Hamiltonians on domains in \( \mathbb{R}^N \), exponential growth of the number of local minima in \( N \) has recently been proved rigorously [Fy].

In the specific models of type IIb flux compactifications on a CY 3-fold \( X \), the number of variables is \( b_3(X) \). As mentioned above, for a typical CY 3-fold, \( b_3 \) is often around 300 and sometimes as high as 1000 (cf. [GHJ, CO]), and therefore the scalar potential \( V_W \) in (11) is a function of this number of variables. By naive counting of variables one would thus arrive at a figure like \( 10^{500} \) for such models. The more sophisticated estimate \( N_{\text{vac}} \approx L^{b_3} / b_3! f(b_3) \) in flux compactifications (see §1.6 for the notation) does not supplant the naive counting argument since the order of magnitude of \( f(b_3) \) is unknown. We recall that \( f(b_3) \) is the integral over \( \mathcal{C} \) of the Gaussian integral in (32) (see (125). The Gaussian integral for \( K_{\text{crit}} \) in that line resembles to some extent the integral formula for the expected number of critical points in spin glass theory, which has exponential growth (see e.g. [Fy]).

Although the naive counting of variables or the analogy to complexity of energy landscapes bring some insight into vacuum counting, we now point out some pitfalls in estimating numbers of vacua or the coefficient \( f(b_3) \) in flux compactifications on this basis.

1. The critical point equation (12) is \( C^\infty \) but not holomorphic, so vacua are critical points of a real system of equations, and it is not obvious how many connection critical points to expect even a polynomial of a given degree to have. This number depends on the connection, and is studied in detail in [DSZ1, DSZ2] and in the present paper.

2. A flux superpotential \( W \) is not a polynomial and it is not clear how to assign it a ‘degree’ which reflects its number of critical points on all of Teichmüller space, or
equivalently, the number of critical points in \( C \) corresponding to the \( \Gamma \)-orbit of \( W \). Examples (e.g. in §6.4.1) show that this number can be relatively small.

(3) It seems reasonable to say that it is the number of fluxes rather than the number of critical points per flux that dominates the number of vacua. In flux compactifications, the landscape should therefore be viewed as the graph of the scalar potential \( V_W(Z) \) on \( C \times S \), i.e. as a function of both variables \( W, Z \), and the local minima should be viewed as pairs \((W_G, Z)\) with \( G \in H^3(X, \mathbb{Z} \oplus \sqrt{-1}\mathbb{Z}) \) and with \( Z \in \text{Crit}(W_G) \).

(4) However (see the problems below) it is not straightforward to define ‘per vacua’, since the tadpole constraint is hyperbolic, and the total number of lattice points in the shell \( 0 < Q[G] < L \) is infinite.

(5) In estimating \( K^\text{crit}(Z) \) we are fixing \( Z \) in the interior of \( C \). But there could exist singular points of \( C \) at which \( K^\text{crit}(Z) \) blows up (see [DD1] for discussion of conifold points). It would also be interesting to study \( K^\text{crit}(Z) \) as \( Z \to \partial C \).

(6) As mentioned in §1.6 (see also §7.3), there may be a significant difference between the order of magnitude of the density of critical points and of the number of critical points, since \( C \) is an incomplete Kähler manifold of possibly quite small volume. See [LS1] for the current state of the art on the volume. There is no analogue of the small volume of the configuration space in spin glass complexity.

(7) The tadpole constraint (1) becomes much more highly constraining as the number \( b_3 \) of variables increases for fixed \( L \) and is responsible for the factor \( \frac{1}{(b_3)} \) in Theorem 1.4. Again, no such feature exists in complexity estimates in spin glasses.

7.2. Problems. The issues mentioned above (and the detailed heuristics in §7.3) suggest a number of problems. The ultimate goal is:

**Problem 7.1.** Does string theory contain a vacuum consistent with the standard model, and if so, how many? Find examples of Calabi-Yau manifolds, and any other postulated structures, for which it is certain that such a vacuum exists.

Now testing consistency with the standard model requires elucidating far more structure of a candidate vacuum – the gauge group, the matter content, and so forth – than we are considering here. To address this ultimate problem, one would need many more statistical results, along the lines set out in [Do]. However one can make arguments (admittedly quite speculative at this point) that the dominant multiplicity in vacuum counting arises from the multiplicity of flux vacua we are discussing here. An important problem in this context is

**Problem 7.2.** How large does \( L \) need to be to ensure that there exists a vacuum with

\[
|W_G(Z)|^2 \leq \lambda_* \tag{123}
\]

for a specified \( \lambda_* \)? In that case, how many such vacua are there? Find examples of Calabi-Yau manifolds where it is certain that such a vacuum exists.

To solve this problem for type IIb flux compactifications, we would need to sharpen Theorem 1.4 in many ways which lead to the subsequent problems stated below.

The constraint (123) on \( |W_G(Z)|^2 \) is a simple example of ‘consistency with the standard model.’ If the real world were (counter-factually) exactly supersymmetric, this would be the constraint that the vacuum should have a cosmological constant \( V_W(Z) = -3|W_G(Z)|^2 \) (as in (11)) consistent with the known value. While the physical discussion
requires taking supersymmetry breaking into account, as discussed in [DD2], vacua can exist in which supersymmetry is broken by effects not taken into account here, making additional contributions to the vacuum energy which lift the exact vacuum energy to be consistent with the known value (essentially, zero). For such a vacuum, the quantity $3|W_G(Z)|^2$ would be the mass squared of the gravitino, a quantity which could be constrained by physical observations.

An independent motivation for (123) is that some proposals for stabilizing the moduli we did not discuss, such as that of [KKLT], are believed only to work under such a constraint.

In any case, as discussed in [DD1] (§3.3), one can count such vacua by choosing the test function to be $\theta(\lambda_\ast - |W_G(Z)|^2)$, where $\theta(x) = 1$ for $x > 0$ and $= 0$ for $x \leq 0$. This test function is not homogeneous but can be handled by the methods of this paper (loc. cit.).

Theorem 1.4 is asymptotic in $L$ and we have also analyzed to some degree the $b_3$ dependence. But as mentioned in §1.6, $L$ depends on the topology of $X$. There, we stated that in many examples $L \simeq Cb_3$ with $1/3 \leq C \leq 3$. To bridge one gap between Theorem 1.4 and Problem 7.2, we state:

**Problem 7.3.** How are the order of magnitudes of $b_3(X)$ and $L$ of (36) related as $X$ varies over topologically distinct Calabi-Yau manifolds?

We have already mentioned the importance of obtaining effective estimates in $b_3$ of the coefficient (24) in Theorem 1.4:

**Problem 7.4.** Obtain an effective estimate of $K_{\text{crit}}(Z)$ and of its integral over $C$ in $b_3$. Also, obtain such an estimate of the remainder.

Among the difficulties with this problem is that $K_{\text{crit}}(Z)$ depends on special features of the moduli space $C$ which depend on more than just the dimension $b_3$ and which may change in an irregular way as the dimension increases. We consider this problem below in §7.3.

To gain insight into the size of the leading coefficient (24), one could write the principal term in Theorem 1.4 in the form $f_{b_3}^{b_3} \times f(b_3)$ that is often used in string theory (cf. §1.6), with $f(b_3)$ the Gaussian integral in (32) (see especially (125)). As mentioned above, it is natural to try to separate out the effects of the number of fluxes and the number of vacua per flux, or more precisely:

1. the number of fluxes $G$ satisfying the tadpole constraint with a critical point in a compact subset $K \subset C$;
2. the number of critical points ‘per flux’, or more precisely per $\Gamma$-orbit of fluxes, in $K$ (see §6.4.1 to clarify this distinction);
3. the total number of critical points in $K$ of all fluxes satisfying the tadpole constraint.

We can define the first quantity precisely as the sum

$$\Theta_K(L) = \sum_{G \in H^3(X, \mathbb{Z} \oplus i\mathbb{Z}) : Q[G] \leq L} \theta \left( \sum_{Z \in C : \nabla W_G(Z) = 0} \chi_K(Z) \right).$$

Thus, the problem we pose is:
Problem 7.5. Determine the asymptotics of $\Theta_K(L)$ as $L \to \infty$.

The second quantity is the ratio $N_K(L)/\Theta_K(L)$. A possibly more tractable way to restate this problem is in terms of the ‘average number of critical points’ of a superpotential $W_G$ in $K$. To define ‘average’ we need to introduce a probability measure on $\mathcal{F}$ which is compatible with $\chi_Q dW$. The most natural probability measures seem to be the normalized Gaussian measures $\gamma_{Z_0}$ on the spaces $S_{Z_0}$ defined by the inner product $Q_{Z_0}$. Thus, we ask for the average number of critical points of $W \in S_{Z_0}$ with respect to $\gamma_{Z_0}$. It would be interesting to study the number of critical points in a fixed $K \subset \mathbb{C}$ or in all of $\mathbb{C}$ or indeed in all of Teichmüller space (which corresponds to counting critical points in $\mathbb{C}$ for a $\Gamma$-orbit of fluxes).

We observe that $W \in S_{Z_0}$ has a critical point at $Z$ if and only if $W \in S_{Z_0} \cap S_{Z}$. In the case of flux superpotentials, $\dim S_{Z_0} = \frac{1}{2} \dim \mathcal{F}$ so for generic pairs $Z, Z_0$, $S_{Z_0} \cap S_{Z} = \{0\}$. Thus, $E_{Z_0}(\#\text{Crits}(W))$ will be an integral over the special variety $\Sigma_{Z_0} = \{Z \in \mathbb{C} : \dim S_{Z_0} \cap S_{Z} > 0\}$. This variety is obviously stratified by $h^{2,1}$ strata $\Sigma_d$ on which the dimension $d$ takes the values $d = 1, 2, \ldots, h^{2,1}$, and $E_{Z_0}(\#\text{Crits}(W))$ is a sum of integrals over each strata.

Problem 7.6. Determine the asymptotics of $E_{Z_0}(\chi_Q Z_0(G/L)\#\text{Crits}(W_G))$.

We also recall that in Theorem 1.4 we ignored the effect of the discriminant variety and the boundary of the region of $\mathcal{C}$.

Problem 7.7. Estimate the remainder if $\psi$ does not vanish near the discriminant variety $D$, or if $\psi$ is a characteristic function of a smooth region $K \subset \mathbb{C}$. Investigate the boundary behavior as $K$ fills out to $\mathbb{C}$.

An analogous problem about studying accumulation of lattice points around boundaries of domains on non-degenerate surfaces is studied in [Ze1].

7.3. Heuristic estimate of the critical point density. We now present a heuristic estimate on the $b_3$-dependence of the critical point density (relative to the Weil-Petersson volume form)

$$K^{\text{crit}}(Z) = \frac{1}{b_3! \sqrt{\det \Lambda_Z}} \int_{\mathcal{H}_Z \oplus \mathbb{C}} \left| \det H^* H - |x|^2 I \right| e^{-(\Lambda_Z^{-1} H, H)_\mathbb{R} - |x|^2} dH dx \tag{124}$$

for $Z$ in regions of moduli space where the norm of $\Lambda_Z$ satisfies bounds independent of $b_3$. We recall (cf. Proposition 6.2) that $\Lambda_Z$ is the Hodge metric, hence we are studying the density of critical points in regions $K \subset \mathcal{C}$ where the absolute values of the eigenvalues of the Ricci curvature of the Weil-Petersson metric $\omega_{WP}$ are bounded by a uniform constant. In the notation $N_{\text{vac}}(L) \sim \frac{L^{b_3}}{b_3!} f(b_3)$, we have

$$f(b_3) = \int_{\mathcal{C}} \chi_K(Z) \frac{1}{\sqrt{\det \Lambda_Z}} \int_{\mathcal{H}_Z \oplus \mathbb{C}} \left| \det H^* H - |x|^2 I \right| e^{-(\Lambda_Z^{-1} H, H)_\mathbb{R} - |x|^2} dH dx, \tag{125}$$

where $\mathcal{K}$ is the region in which we are counting the critical points.
Our heuristic estimate is that the Gaussian integral (i.e. $b_3!K^{\text{crit}}(Z)$) has growth rate $(b_3/2)!N_{\mu}^{b_3}$ for $Z$ in a region $K = K_\mu$ of moduli space where $||\Lambda_Z|| \leq \mu$. Here, $N_{\mu}$ is a constant depending only on $\mu$. It follows that $K^{\text{crit}}(Z)$ would have the decay rate $b_3^{b_3/2}$ for $Z$ in $K_\mu$. We note that this heuristic estimate is consistent with the heuristic estimate given by Ashok-Douglas [AD] that $K^{\text{crit}}(Z)$ should have the same order of magnitude as $\text{Ind}(Z)$ (116). By Proposition 6.3, $b_3!\text{Ind}(Z)$ is a differential form depending polynomially on the curvature. The density of $b_3!\text{Ind}(Z)$ relative to $d\text{Vol}_{WP} = \frac{\omega^b_2}{b_3/2}$ thus has the growth $(b_3/2)!N_{\mu}^{b_3}$ we predict. We present the new heuristic to give evidence that the absolute value only changes the coefficient and not the order of magnitude in vacuum counting.

Before going into the heuristic estimate, we first discuss the consequences for vacuum counting. As mentioned in the introduction, it has been tentatively conjectured at this time of writing (Z. Lu) that the Weil-Petersson volume of $K_\mu$ is bounded above by the volume of a ball of radius $r(\mu)$ in $\mathbb{C}^{b_3/2}$ depending only on $\mu$, and the latter volume decays like $\frac{1}{(b_3/2)!}$. Thus it would appear that $N_{\text{vac},K_\mu}(L) \sim \frac{(C_1 LN_{\mu})^{b_3}}{b_3!}$. We include a constant $C_1$ to take into account the dependence on various parameters including $r(\mu)$, factors of $\pi$ and so on. If we then take the (often) observed value $L \sim C b_3$ with $C \in [\frac{1}{3}, 3]$, then the number of vacua in $K_\mu$ satisfying the tadpole constraint would grow at an exponential rate in $b_3$.

We now explain the heuristic estimate regarding the order of magnitude of $K^{\text{crit}}(Z)$ (24): the latter depends on two inputs, the subspace $\mathcal{H}_Z$ (or equivalently the orthogonal projection $P_Z$ onto $\mathcal{H}_Z$) and the eigenvalues of $\Lambda_Z$. To obtain upper and lower bounds on $K^{\text{crit}}(Z)$ we note that

$$2P_Z \leq \Lambda_Z \leq \mu_{\text{max}}(Z)P_Z,$$

(126)

where $\mu_{\text{max}}(Z)$ is the maximum eigenvalue of $\Lambda_Z$. We recall here that $\Lambda_Z$ is the matrix of the Hodge metric (see (30)), and its eigenvalues can be estimated in terms of the Weil-Petersson metric and its curvature (cf. [Lu]). In particular, its minimum eigenvalue satisfies $\mu_{\text{min}}(Z) \geq 2$, and that explains the lower bound $2P_Z$ in (126). For most CY 3-folds $X$, the Weil-Petersson metric on $\mathcal{C}$ is incomplete, and $\mu_{\text{max}}(Z) \to \infty$ as $Z$ tends to the boundary (Z. Lu).

By (126), we have

$$J_-(\mu, P_Z) \leq (b_3!)K^{\text{crit}}(Z) \leq J_+(\mu, P_Z), \quad (\forall \mu \geq \mu_{\text{max}}(Z)),$$

(127)

where

$$J_+(\mu, P_Z) := \frac{1}{2^{b_3/2-1}} \int_{\mathcal{H}_Z \otimes \mathbb{C}} \left| \det H^*H - |x|^2 I \right| e^{-(\mu^{-1}\text{Tr}H^*H - |x|^2)} dH dx,$$

(128)

and

$$J_-(\mu, P_Z) := \frac{1}{\mu^{(b_3/2-1)}} \int_{\mathcal{H}_Z \otimes \mathbb{C}} \left| \det H^*H - |x|^2 I \right| e^{-(2^{-1}\text{Tr}H^*H - |x|^2)} dH dx,$$

(129)

Thus we obtain upper and lower bounds for the density in regions $K_\mu \subset \mathcal{C}$ for which the absolute values of the eigenvalues of the Hodge metric relative to the Weil-Petersson metric...
metric satisfy \( \mu_{\text{max}}(Z) \leq \mu \). We have bounded the determinant of \( \Lambda \) by a power of an extremal eigenvalue, but it could also be identified with the volume density of the Hodge metric. We note that the lower bound tends to zero and the upper bound tends to infinity in \( \sim \pm b_3 \) powers of \( \mu_{\text{max}}(Z) \) as \( Z \rightarrow \partial \mathcal{C} \) when the Weil-Petersson metric is incomplete and the norm of the Ricci curvature of \( \omega_{WP} \) tends to infinity.

We now estimate \( J_{\pm}(\mu, P_Z) \) under the assumption that \( \mathcal{H}_Z \) is a ‘sufficiently random’ subspace. The subspace \( \mathcal{H}_Z \) is a real subspace of dimension \( b_3 - 2 \) of \( \text{Sym}(b_3/2 - 1, \mathbb{C}) \), but by modifying the definition of the complex structure it becomes a complex \( b_3/2 \)-dimensional one. Hence, we may view \( Z \rightarrow \mathcal{H}_Z \) as a map \( \mathcal{C} \rightarrow \text{Gr}(b_3/2 - 1, \text{Sym}(b_3/2 - 1, \mathbb{C})) \) to the complex Grassmannian of \( b_3/2 - 1 \) dimensional complex subspaces. Lacking knowledge of the distribution of the image of \( Z \rightarrow \mathcal{H}_Z \), we make the assumption that it is random, or more precisely we approximate \( J_{\pm}(\mu, P_Z) \) by the expected value of \( J_{\pm}(\mu, P) \), where \( P \) is the projection corresponding to a random element \( \mathcal{H} \in \text{Gr}(b_3/2 - 1, \text{Sym}(b_3/2 - 1, \mathbb{C})) \).

This approximation by the expected value seems to be reasonable because Grassmannians \( \text{Gr}(k, N) \) are examples of Gromov-Milman ‘Levy families’ of Riemannian manifolds for which concentration of measure phenomena hold as \( N \rightarrow \infty \) [GM, Ta]. Concentration of measure refers to a metric space \((X, d)\) with a probability measure \( P \) and a concentration function \( \alpha(P, t) \), which is the smallest number such that the measure of a set \( A \) and the metric tube \( A_t = \{ x : d(x, A) < t \} \) around \( A \) are related by \( P(A) \geq 1/2 \implies P(A_t) \geq 1 - \alpha(P, t) \). If \( f \) is a Lipschitz function and if \( M_f \) is a median for \( f \), we put \( A = \{ x : f(x) \leq M_f \} \), and then \( P(\{|f - M_f| > t\}) \leq 2\alpha(P, \frac{t}{\|f - M_f\|_{Lip}}) \).

Concentration of measure occurs if \( \alpha(P, t) \) decays rapidly in \( t \), and thus \( f \) is highly concentrated around its median. In a Lévy family \((X_N, d_N)\), the functions \( \alpha N(P, t) \) decay at ever faster rates depending on \( N \). For instance on the unit \( N \)-sphere \( S^N \), the rate is (a universal constant times) \( e^{-\frac{(N-1)^2}{2} t^2} \).

In our setting, the family consists of Grassmannians \( \text{Gr}(b_3/2 - 1, \text{Sym}(b_3/2 - 1, \mathbb{C})) \) equipped with the invariant probability measure \( d\nu \) and with the standard bi-invariant metric. It is pointed out in [GM] that \( \text{Gr}(k, N) \) is a Lévy family for fixed \( k \) (see Sect. (3.3) of [GM]), and the same argument should apply to \( k_N \sim N/2 \). Moreover, \( \{U(N)\} \) with its Haar probability measure and bi-invariant metric is Lévy, and by Sect. (2.1) of [GM] its quotients should be. The function \( f \) is \( J_{\pm}(\mu, P) \) for fixed \( \mu \). Since we are mainly interested in factorial dependencies, we set \( \mu = 1 \) and change the exponent \( 2^{-1} \) to 1 to make the Gaussian a probability measure. In general, the result would be modified by a \( \pm b_3 \) power of \( \mu \). In this heuristic discussion, we will not attempt to determine \( \alpha N(P, t) \) or \( M_f \) but will assume that \( \alpha(P, \frac{t}{\|f - M_f\|_{Lip}}) \) has rapid decrease in \( t \) which improves with the dimension. We also note that when \( \alpha(P, t) \) is small, we can replace the median of \( J_{\pm}(\mu, P) \) (with \( \mu = 1 \)) by its mean

\[
\int_{\text{Gr}(b_3/2 - 1, \text{Sym}(b_3/2 - 1, \mathbb{C}))} \int_{\mathcal{H} \oplus \mathbb{C}} |\det(H^*H - |x|^2 I)|e^{-TrH^*H - |x|^2}dHdx \, d\nu(\mathcal{H})
\]

with a small error (cf. [Ta]). This mean equals

\[
\int_{\text{Sym}(b_3/2 - 1, \mathbb{C}) \oplus \mathbb{C}} |\det(H^*H - |x|^2 I)|e^{-TrH^*H - |x|^2}dHdx \quad (130)
\]

since both measures are invariant probability measures and are therefore equal. Here we ignore factors of \((2\pi)\) (etc.) for the sake of simplicity, since we are primarily interested in the factorially growing quantities. Due to the concentration of measure, the spaces
\( H_Z \) would have to be very ‘rare events’ if \( J_{\pm}(\mu, P_Z) \) differed appreciably from its mean. We note that since \( H^3_Z \) is a complex polarization, \( P_Z \) has special features that do not hold for random subspaces, but we have no reason to believe that these special features bias \( J(\mu, P_Z) \) away from its mean.

We now observe that (130) (with any choice of \( \mu \)) is similar to the integral for the density of critical points for holomorphic sections of \( \mathcal{O}(N) \rightarrow \mathbb{CP}^m \) with \( m = b_3/2 - 1 \) with respect to the Fubini-Study connection for a fixed degree \( N \) [DSZ2] (§4). There, the \( \Lambda_Z \) matrix was (for every \( Z \)) a two-block diagonal matrix with a large scalar block and a \( 1 \times 1 \) scalar block. When \( \mu = 1 \) (130) agrees with that \( \mathcal{O}(N) \rightarrow \mathbb{CP}^m \) density in the case \( N = 1 \). As noted in [DSZ2], the total number of critical points of a given Morse index appears to grow at a rate \( N^m \) times a rational quantity in \( m \) as \( m \rightarrow \infty \). This growth rate may also be easily verified for the Euler characteristic \( c_m(T^{*1,0} \otimes \mathcal{O}(N)) \), i.e. the alternating sum over the Morse indices, which is given by

\[
c_m(T^{*1,0} \otimes \mathcal{O}(N)) = \left( \frac{c(\mathcal{O}(N) - 1)^{m+1}}{c(\mathcal{O}(N))}, \mathbb{CP}^m \right) = \frac{N - 1)^{m+1} + (-1)^m}{N}.\]

Since the volume of \( \mathbb{CP}^m \) is \( \frac{1}{m!} \) this would imply that the density of critical points grows like \( m! \) with the dimension. On this basis, we would expect that \( J_{\pm}(\mu, P_Z) \) for \( \mu \simeq 1 \) grows with the dimension at the rate \((b_3/2)!N_{b_3}^{b_3}\) for some \( N_{b_3} > 0 \).

We note that the Ashok-Douglas heuristic that the density of critical points should have the same order of magnitude as the index density is indeed correct in the setting of \( \mathcal{O}(N) \rightarrow \mathbb{CP}^m \). Further, the origin of the factorials \((b_3/2)!\) is essentially in both the \( C \) and \( \mathbb{CP}^m \) settings.

Thus our heuristics give \( K^{\text{crit}}(Z) \sim \frac{(b_3/2)!N_{b_3}^{b_3}}{b_3!} \). If we integrate over \( K_\mu \) and apply the conjectural volume bound \( \frac{1}{(b_3/2)!} \) for \( K_\mu \), we would get roughly \( \frac{L_{b_3}^3 N_{b_3}^{b_3}}{b_3!} \). Further applying the observed relation \( L \sim Cb_3 \) with \( C \in [1/3, 3] \) gives an exponential growth rate for numbers of vacua in \( K_\mu \).

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