

INTEGRATING PRODUCTS OF QUADRATIC FORMS

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ABSTRACT. We prove that if $q_1, \dots, q_m : \mathbb{R}^n \rightarrow \mathbb{R}$ are quadratic forms in variables x_1, \dots, x_n such that each q_k depends on at most r variables and each q_k has common variables with at most r other forms, then the average value of the product $(1 + q_1) \cdots (1 + q_m)$ with respect to the standard Gaussian measure in \mathbb{R}^n can be approximated within relative error $\epsilon > 0$ in quasi-polynomial $n^{O(1)} m^{O(\ln m - \ln \epsilon)}$ time, provided $|q_k(x)| \leq \gamma \|x\|^2 / r$ for some absolute constant $\gamma > 0$ and $k = 1, \dots, m$. When q_k are interpreted as pairwise squared distances for configurations of points in Euclidean space, the average can be interpreted as the partition function of systems of particles with mollified logarithmic potentials. We sketch a possible application to testing the feasibility of systems of real quadratic equations.

1. INTRODUCTION AND MAIN RESULTS

Integration of high degree multivariate polynomials is computationally difficult and no efficient algorithms are known except in few special cases, when the polynomials have a rather simple algebraic structure (for example, close to a power of a linear form on a simplex, cf. [B+11]), or have some very nice analytic properties (slowly varying or, most notably, log-concave, cf. [LV07]). Since a general n -variate polynomial p of degree d is defined by $\binom{n+d}{d}$ parameters (for example, coefficients), the problem becomes interesting for large n and d only if p has some special structure (such as the product of low-degree polynomials), which allows us to define p using much fewer parameters.

In this paper, we integrate products of quadratic forms with respect to the Gaussian measure in \mathbb{R}^n . We relate the problem to partition functions of mollified logarithmic potentials and to testing the feasibility of systems of real quadratic equations.

Our algorithms are deterministic and based on the method of polynomial interpolation, which has been recently applied to a variety of partition functions in

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combinatorial (discrete) problems, cf. [Ba16]. In continuous setting, the method was applied to computing partition functions arising in quantum models [B+19], [H+19].

(1.1) Quadratic forms on \mathbb{R}^n . We consider Euclidean space \mathbb{R}^n endowed with the standard inner product

$$\langle x, y \rangle = x_1 y_1 + \dots + x_n y_n \quad \text{for } x = (x_1, \dots, x_n) \quad \text{and } y = (y_1, \dots, y_n)$$

and corresponding Euclidean norm

$$\|x\| = \sqrt{\langle x, x \rangle} = \sqrt{x_1^2 + \dots + x_n^2} \quad \text{for } x = (x_1, \dots, x_n).$$

Let $q_1, \dots, q_m : \mathbb{R}^n \rightarrow \mathbb{R}$ be quadratic forms defined by

$$(1.1.1) \quad q_k(x) = \frac{1}{2} \langle Q_k x, x \rangle \quad \text{for } k = 1, \dots, m,$$

where Q_1, \dots, Q_m are $n \times n$ real symmetric matrices.

Our first result concerns computing the integral

$$(1.1.2) \quad \frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} (1 + q_1(x)) \cdots (1 + q_m(x)) e^{-\|x\|^2/2} dx.$$

The idea of the interpolation method is to consider (1.1.2) as a one-parameter perturbation a much simpler integral, in our case, of

$$(1.1.3) \quad \frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} e^{-\|x\|^2/2} dx = 1.$$

For the method to work, one should show that there are no zeros in the vicinity of a path in the complex plane which connects (1.1.2) and (1.1.3). We prove the following result.

(1.2) Theorem. *There is an absolute constant $\gamma > 0$ (one can choose $\gamma = 0.151$) such that the following holds. Let $q_k : \mathbb{R}^n \rightarrow \mathbb{R}$, $k = 1, \dots, m$, be quadratic forms. Then*

$$\frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} (1 + \omega q_1(x)) \cdots (1 + \omega q_m(x)) e^{-\|x\|^2/2} dx \neq 0$$

for all $\omega \in \mathbb{C}$ such that $|\omega| \leq \gamma$, provided

$$|q_k(x)| \leq \frac{1}{\max\{m, n\}} \|x\|^2 \quad \text{for } k = 1, \dots, m.$$

By interpolation, for any constant $0 < \gamma' < \gamma$, fixed in advance, we obtain an algorithm which, given quadratic forms $q_1, \dots, q_m : \mathbb{R}^n \rightarrow \mathbb{R}$, computes (1.1.2) within relative error $0 < \epsilon < 1$ in quasi-polynomial $n^{O(1)} m^{O(\ln m - \ln \epsilon)}$ time provided

$$(1.2.1) \quad |q_k(x)| \leq \frac{\gamma'}{\max\{m, n\}} \|x\|^2 \quad \text{for } k = 1, \dots, m.$$

Note that by Theorem 1.2 and (1.1.3), the value of (1.1.2) is positive, as long as (1.2.1) holds.

Some remarks are in order.

First, we note that the integrand in (1.1.2) can vary wildly. Indeed, for large n the bulk of the standard Gaussian measure in \mathbb{R}^n is concentrated in the vicinity of the sphere $\|x\| = \sqrt{n}$, see for example, Section V.5 of [Ba02]. Assuming that $m = n$, we can choose $q_k(x) \sim \|x\|^2/n$ so that (1.2.1) is satisfied. Then, in the vicinity of the sphere $\|x\| = \sqrt{n}$, the product $(1 + q_1(x)) \cdots (1 + q_n(x))$ in (1.1.2) varies within an exponential in n factor, and is not at all well-concentrated.

Second, if the quadratic forms q_1, \dots, q_m exhibit a simpler combinatorics, we can improve the bounds accordingly. We prove the following result.

(1.3) Theorem. *There is an absolute constant $\gamma > 0$ (one can choose $\gamma = 0.151$) such that the following holds. Let $q_k : \mathbb{R}^n \rightarrow \mathbb{R}$, $k = 1, \dots, m$, be quadratic forms. Suppose further that each form depends on not more than r variables among x_1, \dots, x_n and that each form has common variables with not more than r other forms. Then*

$$\frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} (1 + \omega q_1(x)) \cdots (1 + \omega q_m(x)) e^{-\|x\|^2/2} dx \neq 0$$

for all $w \in \mathbb{C}$ such that $|\omega| \leq \gamma$, provided

$$|q_k(x)| \leq \frac{\|x\|^2}{r} \quad \text{for } k = 1, \dots, m.$$

By interpolation, for any constant $0 < \gamma' < \gamma$, fixed in advance, we obtain an algorithm which, given quadratic forms $q_1, \dots, q_m : \mathbb{R}^n \rightarrow \mathbb{R}$ as in Theorem 1.3 computes (1.1.2) within relative error $0 < \epsilon < 1$ in quasi-polynomial $n^{O(1)} m^{\ln m - \ln \epsilon}$ time provided

$$(1.3.1) \quad |q_k(x)| \leq \frac{\gamma'}{r} \|x\|^2 \quad \text{for } k = 1, \dots, m.$$

We prove Theorems 1.2 and 1.3 in Section 3 and describe the algorithm for computing (1.1.2) in Section 4. In Section 2, we discuss connections with systems of particles with mollified logarithmic potentials and possible applications to testing the feasibility of systems of multivariate real quadratic equations.

2. CONNECTIONS AND POSSIBLE APPLICATIONS

(2.1) Partition functions of mollified logarithmic potentials. Let $n = ds$ and let us interpret $\mathbb{R}^n = \mathbb{R}^d \oplus \cdots \oplus \mathbb{R}^d$ as the space of all ordered s -tuples (v_1, \dots, v_s) of points $v_i \in \mathbb{R}^d$. Hence the distance between v_i and v_j is $\|v_i - v_j\|$.

Let us fix some set E of m pairs $\{i, j\}$ of indices $1 \leq i < j \leq s$ and suppose that the energy of a set of points (v_1, \dots, v_s) is defined by

$$(2.1.1) \quad - \sum_{\{i,j\} \in E} \ln(1 + \alpha \|v_i - v_j\|^2) + \frac{1}{2} \sum_{i=1}^n \|v_i\|^2,$$

where $\alpha > 0$ is a parameter. The first sum in (2.1.1) indicates that there a repulsive force between any pair $\{v_i, v_j\}$ with $\{i, j\} \in E$ (so that the energy decreases if the distance between v_i and v_j increases), while the second sum indicates that there is a force pushing the points towards 0 (so that the energy decreases when each v_i approaches 0). When $\alpha = 0$, the repulsive force disappears altogether, and $\alpha \rightarrow +\infty$, the repulsive force behaves as a Coulomb's force with logarithmic potential, since

$$\lim_{\alpha \rightarrow +\infty} \ln(1 + \alpha \|v_i - v_j\|^2) - \ln \alpha = 2 \ln \|v_i - v_j\|.$$

Thus the integral

$$(2.1.2) \quad \frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} \prod_{\{i,j\} \in E} (1 + \alpha \|v_i - v_j\|^2) e^{-(\|v_1\|^2 + \dots + \|v_s\|^2)/2} dx,$$

which is a particular case of (1.1.2), can be interpreted as the partition function of points with “mollified” or “damped” logarithmic potentials. One can think of (2.1.2) as the partition function for particles with genuine logarithmic potentials, provided each particle is confined to its own copy of \mathbb{R}^d among a family of parallel d -dimensional affine subspaces in some higher-dimensional Euclidean space.

The integral (2.1.2) can be considered as a ramification of classical Selberg-type integrals for logarithmic potentials:

$$(2.1.3) \quad \begin{aligned} & \frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} \prod_{1 \leq i < j \leq n} |x_i - x_j|^{2\gamma} e^{-(x_1^2 + \dots + x_n^2)/2} dx_1 \cdots dx_n \\ & = \prod_{j=1}^n \frac{\Gamma(1 + j\gamma)}{\Gamma(1 + \gamma)}, \end{aligned}$$

see for example, Chapter 17 of [Me04]. The integral (2.1.3) corresponds to points in \mathbb{R}^1 and a similar integral is computed explicitly for points in \mathbb{R}^2 (and $\gamma = 1$), see Section 17.11 of [Me04]. For higher dimensions d no explicit formulas appear to be known.

In contrast, we compute integrals (2.1.2) approximately for certain values of α , but we allow arbitrary dimensions and can choose an arbitrary set of pairs of interacting points (and we can even choose different α s for different pairs of points). Theorem 1.3 can be interpreted as the absence of phase transition in the Lee - Yang sense [YL52], if α is sufficiently small. For example, if the set E consists of all $\binom{s}{2}$ pairs $\{i, j\}$, Theorem 1.3 implies that there is no phase transition (and the integral can be efficiently approximated) if

$$\alpha < \frac{\beta}{\max\{d, s\}}$$

for some absolute constant $\beta > 0$.

(2.2) Applications to systems of quadratic equations. Every system of real polynomial equations can be reduced to a system of quadratic equations, as one can successively reduce the degree by introducing new variables via substitutions of the type $z := xy$. A system of quadratic equations can be solved in polynomial time when the number of equations is fixed in advance, [Ba93], [GP05], but as the number of equations grows, the problem becomes computationally hard.

Here we are interested in the systems of equations of the type

$$(2.2.1) \quad q_k(x) = 1 \quad \text{for } k = 1, \dots, m,$$

where $q_k : \mathbb{R}^n \rightarrow \mathbb{R}$ are positive semidefinite quadratic forms. Such systems naturally arise in problems of distance geometry, where we are interested to find out if there are configurations of points in \mathbb{R}^d with prescribed distances between some pairs of points and in which case q_k are scaled squared distances between points, see [CH88], [L+14] and Section 2.1. Besides, finding if a system of homogeneous quadratic equations has a non-trivial solution

$$(2.2.2) \quad q_k(x) = 0 \quad \text{for } k = 1, \dots, m \quad \text{and} \quad \|x\| = 1$$

can be reduced to (2.2.1) with positive definite forms q_k by adding $\|x\|^2$ to the appropriately scaled equations in (2.2.2).

Suppose that

$$(2.2.3) \quad \sum_{k=1}^m q_k(x) = \frac{\|x\|^2}{2}$$

in (2.2.1). By itself, the condition (2.2.3) is not particularly restrictive: if the sum in the left hand side of (2.2.3) is positive definite, it can be brought to the right hand side by an invertible linear transformation of x .

Let us choose an $\alpha > 0$ such that the scaled forms αq_k satisfy (1.3.1), so that the integral

$$(2.2.4) \quad \frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} (1 + \alpha q_1(x)) \cdots (1 + \alpha q_m(x)) e^{-\|x\|^2/2} dx$$

can be efficiently approximated. We would like to argue that the value of the integral (2.2.4) can provide a reasonable certificate which allows one to distinguish systems (2.2.1) with many “near solutions” x from the systems that are far from having a solution.

We observe that the system (2.2.1) has a solution if and only if the system

$$(2.2.5) \quad q_k(x) = t \quad \text{for } k = 1, \dots, m$$

has a solution $x \in \mathbb{R}^n$ for any $t > 0$.

Let us find $0 < \beta < 1$ such that

$$(2.2.6) \quad 2m \left(\frac{1}{\beta} - \frac{1}{\alpha} \right) = \frac{n}{1 - \beta}.$$

Indeed, the equation (2.2.6) always has a (necessarily unique) solution $0 < \beta < 1$, since for $\beta \approx 0$ the left hand side is bigger than the right hand side, while for $\beta \approx 1$ the right hand side is bigger than the left hand side.

Because of (2.2.3), we can rewrite (2.2.4) as

$$(2.2.7) \quad \frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} e^{-\frac{(1-\beta)\|x\|^2}{2}} \prod_{k=1}^m (1 + \alpha q_k(x)) e^{-\beta q_k(x)} dx.$$

We observe that if $\alpha > \beta$ then the maximum value of

$$(1 + \alpha t) e^{-\beta t} \quad \text{for } t \geq 0$$

is attained at

$$(2.2.8) \quad t = \frac{1}{\beta} - \frac{1}{\alpha} > 0$$

and is equal to

$$\frac{\alpha}{\beta} \exp \left\{ \frac{\beta}{\alpha} - 1 \right\} > 1$$

and hence the maximum value of the product of the m factors in (2.2.7) is

$$\left(\frac{\alpha}{\beta} \exp \left\{ \frac{\beta}{\alpha} - 1 \right\} \right)^m$$

and attained if and only if the system (2.2.1) and hence (2.2.6) has a solution x .

Also, if x is a solution to (2.2.5), by (2.2.3), (2.2.6) and (2.2.8), we have

$$\|x\|^2 = 2tm = \frac{n}{1 - \beta}.$$

The Gaussian probability measure in \mathbb{R}^n with density

$$\frac{(1-\beta)^{n/2}}{(2\pi)^{n/2}} e^{-\frac{(1-\beta)\|x\|^2}{2}},$$

is concentrated in the vicinity of the sphere $\|x\|^2 = n/(1-\beta)$, cf., for example, Section V.5 of [Ba05] for some estimates. Therefore, if for the system (2.2.1) there are sufficiently many “near solutions” x , we should have the value of the integral (2.2.4) sufficiently close to

$$\left(\frac{\alpha}{\beta} \exp\left\{\frac{\beta}{\alpha} - 1\right\}\right)^m (1-\beta)^{-\frac{n}{2}},$$

while if the system (2.2.1) is far from having a solution, the value of the integral will be essentially smaller.

3. PROOFS OF THEOREMS 1.2 AND 1.3

Choosing $r = \max\{m, n\}$, we obtain Theorem 1.2 as a particular case of Theorem 1.3. Hence we prove Theorem 1.3 only.

For a real symmetric $n \times n$ matrix Q we denote

$$\|Q\| = \max_{\|x\|=1} \|Qx\|$$

its operator norm.

We start with a simple formula, cf. also [Ba93].

(3.1) Lemma. *Let $q_1, \dots, q_m : \mathbb{R}^n \rightarrow \mathbb{R}$ be quadratic forms,*

$$q_k(x) = \frac{1}{2} \langle Q_k x, x \rangle \quad \text{for } k = 1, \dots, m,$$

where Q_1, \dots, Q_m are $n \times n$ real symmetric matrices such that

$$\sum_{k=1}^m \|Q_k\| < 1.$$

Then

$$(3.1.1) \quad \det^{-\frac{1}{2}} \left(I - \sum_{k=1}^m z_k Q_k \right) = \sum_{k_1, \dots, k_m \geq 0} \frac{z_1^{k_1} \cdots z_m^{k_m}}{k_1! \cdots k_m!} \\ \times \frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} q_1^{k_1}(x) \cdots q_m^{k_m}(x) e^{-\|x\|^2/2} dx,$$

for all $z_1, \dots, z_m \in \mathbb{C}$ such that

$$(3.1.2) \quad |z_1|, \dots, |z_m| \leq 1.$$

Here we take the principal branch of $\det^{-\frac{1}{2}}$ in the left hand side of (3.1.1), which is equal to 1 when $z_1 = \dots = z_m = 0$. The series in the right hand side converges absolutely and uniformly on the polydisc (3.1.2).

Proof. For $z = (z_1, \dots, z_m)$, let

$$Q_z = I - \sum_{k=1}^m z_k Q_k$$

and let

$$q_z(x) = \frac{1}{2} \langle Q_z x, x \rangle = \frac{\|x\|^2}{2} - \sum_{k=1}^m z_k q_k(x).$$

If z_1, \dots, z_m are real and satisfy (3.1.2), then $q_z : \mathbb{R}^n \rightarrow \mathbb{R}$ is a positive definite quadratic form, and, as is well known,

$$\frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} e^{-q_z(x)} dx = \frac{1}{\sqrt{\det Q_z}}.$$

Since both sides of the above identity are analytic in the domain (3.1.2), we obtain

$$\det^{-\frac{1}{2}} \left(I - \sum_{k=1}^m z_k Q_k \right) = \frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} \exp \left\{ -\frac{\|x\|^2}{2} + \sum_{k=1}^m z_k q_k(x) \right\} dx.$$

Expanding the integral in the right hand side into the series in z_1, \dots, z_m , we complete the proof. \square

Next, we extract the integral (1.1.2) from the generating function of Lemma 3.1. Let

$$\mathbb{S}^1 = \{z \in \mathbb{C} : |z| = 1\}$$

be the unit circle and let

$$\mathbb{T}^m = \underbrace{\mathbb{S}^1 \times \dots \times \mathbb{S}^1}_{m \text{ times}}$$

be the m -dimensional torus endowed with the uniform (Haar) probability measure

$$\mu = \mu_1 \times \dots \times \mu_m,$$

where μ_k is the uniform probability measure on the k -th copy of \mathbb{S}^1 . If $s \in \mathbb{Z}^m$, $s = (s_1, \dots, s_m)$, then for the Laurent monomial

$$\mathbf{z}^s = z_1^{s_1} \dots z_m^{s_m},$$

we have

$$\int_{\mathbb{T}^m} \mathbf{z}^s d\mu = \begin{cases} 1 & \text{if } s = 0 \\ 0 & \text{if } s \neq 0. \end{cases}$$

(3.2) Lemma. Let $q_1, \dots, q_m : \mathbb{R}^n \rightarrow \mathbb{R}$ be quadratic forms,

$$q_k(x) = \frac{1}{2} \langle Q_k x, x \rangle \quad \text{for } k = 1, \dots, m,$$

where Q_1, \dots, Q_m are $n \times n$ real symmetric matrices such that

$$\sum_{k=1}^m \|Q_k\| < 1.$$

Then for every $\omega \in \mathbb{C}$ such that $|\omega| < 1$ we have

$$\begin{aligned} & \frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} (1 + \omega q_1(x)) \cdots (1 + \omega q_m(x)) e^{-\|x\|^2/2} dx \\ &= \int_{\mathbb{T}^m} \prod_{k=1}^m (1 + \omega z_k^{-1}) \prod_{(k_1, \dots, k_s)} \left(1 + \frac{1}{2^s} \text{trace}(Q_{k_1} \cdots Q_{k_s}) z_{k_1} \cdots z_{k_s} \right) d\mu, \end{aligned}$$

where the second product is taken over all non-empty ordered tuples (k_1, \dots, k_s) of distinct indices from $\{1, \dots, m\}$.

Proof. From Lemma 3.1, we have

$$\begin{aligned} & \frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} (1 + \omega q_1(x)) \cdots (1 + \omega q_m(x)) e^{-\|x\|^2/2} dx \\ (3.2.1) \quad &= \int_{\mathbb{T}^m} \prod_{k=1}^m (1 + \omega z_k^{-1}) \det^{-\frac{1}{2}} \left(I - \sum_{k=1}^m z_k Q_k \right) d\mu. \end{aligned}$$

Next, we write

$$\begin{aligned} & \det^{-\frac{1}{2}} \left(I - \sum_{k=1}^m z_k Q_k \right) = \exp \left\{ -\frac{1}{2} \ln \det \left(I - \sum_{k=1}^m z_k Q_k \right) \right\} \\ &= \exp \left\{ -\frac{1}{2} \text{trace} \ln \left(I - \sum_{k=1}^m z_k Q_k \right) \right\} = \exp \left\{ \frac{1}{2} \sum_{s=1}^{\infty} \frac{1}{s} \text{trace} \left(\sum_{k=1}^m z_k Q_k \right)^s \right\} \\ &= \exp \left\{ \sum_{s=1}^{\infty} \frac{1}{2^s} \sum_{1 \leq k_1, \dots, k_s \leq m} \text{trace}(Q_{k_1} \cdots Q_{k_s}) z_{k_1} \cdots z_{k_s} \right\} \\ &= \prod_{s=1}^{\infty} \prod_{1 \leq k_1, \dots, k_s \leq m} \exp \left\{ \frac{1}{2^s} \text{trace}(Q_{k_1} \cdots Q_{k_s}) z_{k_1} \cdots z_{k_s} \right\}, \end{aligned}$$

where the product converges absolutely and uniformly on \mathbb{T}^m .

We expand each of the exponential functions into the Taylor series and observe that only square-free monomials in z_1, \dots, z_m contribute to the integral (3.2.1), from which it follows that

$$\begin{aligned} & \frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} (1 + \omega q_1(x)) \cdots (1 + \omega q_m(x)) e^{-\|x\|^2/2} dx \\ &= \int_{\mathbb{T}^m} \prod_{k=1}^m (1 + \omega z_k^{-1}) \prod_{(k_1, \dots, k_s)} \left(1 + \frac{1}{2s} \text{trace}(Q_{k_1} \cdots Q_{k_s}) z_{k_1} \cdots z_{k_s} \right) d\mu, \end{aligned}$$

where the second product is taken over all non-empty ordered tuples of distinct indices $k_1, \dots, k_s \in \{1, \dots, m\}$. \square

Our next goal is to write the integral in Lemma 3.2 as the value of the independence polynomial of an appropriate (large) graph.

(3.3) Independent sets in weighted graphs. Let $G = (V, E)$ be a finite undirected graph with set V of vertices, set E of edges and without loops or multiple edges. A set $S \subset V$ of vertices is called *independent*, if no two vertices from S span an edge of G . We agree that $S = \emptyset$ is an independent set.

Let $w : V \rightarrow \mathbb{C}$ be a function assigning to each vertex a complex *weight* $w(v)$. We define the *independence polynomial* of G by

$$\text{Ind}_G(w) = \sum_{\substack{S \subset V \\ S \text{ independent}}} \prod_{v \in S} w(v).$$

Hence $\text{Ind}_G(w)$ is a multivariate polynomial in complex variables $w(v)$ with constant term 1, corresponding to $S = \emptyset$.

(3.4) Corollary. Let $q_1, \dots, q_m : \mathbb{R}^n \rightarrow \mathbb{R}$ be quadratic forms,

$$q_k = \frac{1}{2} \langle Q_k x, x \rangle \quad \text{for } k = 1, \dots, m,$$

where Q_k are real symmetric $n \times n$ matrices and let $\omega \in \mathbb{C}$ be a complex number.

We define a weighted graph $G = (V, E; w)$ as follows. The vertices of G are all non-empty ordered tuples (k_1, \dots, k_s) of distinct indices $k_1, \dots, k_s \in \{1, \dots, m\}$ and two vertices span an edge of G if they have at least one common index k , in arbitrary positions. We define the weight of the vertex (k_1, \dots, k_s) by

$$\frac{\omega^s}{2s} \text{trace}(Q_{k_1} \cdots Q_{k_s}).$$

Then

$$(3.4.1) \quad \frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} (1 + \omega q_1(x)) \cdots (1 + \omega q_m(x)) e^{-\|x\|^2/2} dx = \text{Ind}_G(w).$$

Proof. From Lemma 3.2 it follows that (3.4.1) holds provided $|\omega|$ and $\|Q_k\|$ for $k = 1, \dots, m$ are small enough. Since both sides of (3.4.1) are polynomials in Q_1, \dots, Q_k and ω , the proof follows. \square

The following criterion provides a sufficient condition for $\text{Ind}_G(w) \neq 0$ for an arbitrary weighted graph G . The result is known as the Dobrushin criterion and also as the Kotecký - Preiss condition for the cluster expansion, see, for example, Chapter 5 of [FV18].

(3.5) Lemma. *Given a graph $G = (V, E)$ and a vertex $v \in V$, we define its neighborhood $\mathcal{N}_v \subset V$ by*

$$\mathcal{N}_v = \{v\} \cup \{u \in V : \{u, v\} \in E\}.$$

Let $w : V \rightarrow \mathbb{C}$ be an assignment of complex weights to the vertices of G . Suppose that there is a function $\rho : V \rightarrow \mathbb{R}_+$ with positive real values such that for every vertex $v \in V$, we have

$$\sum_{u \in \mathcal{N}_v} |w(u)| e^{\rho(u)} \leq \rho(v).$$

Then

$$\text{Ind}_G(w) \neq 0.$$

Proof. See, for example, Section 5.2 of [CF16] for a concise exposition. \square

Now we are ready to prove Theorem 1.3.

(3.6) Proof of Theorem 1.3. Let Q_1, \dots, Q_m be the matrices of the quadratic forms q_1, \dots, q_m , so that

$$q_k(x) = \frac{1}{2} \langle Q_k, x \rangle \quad \text{and} \quad \|Q_k\| \leq \frac{2}{r} \quad \text{for } k = 1, \dots, m.$$

Since each quadratic form q_k depends of at most r variables, we have

$$\text{rank } Q_k \leq r \quad \text{for } k = 1, \dots, m.$$

In particular,

$$(3.6.1) \quad |\text{trace}(Q_{k_1} \cdots Q_{k_s})| \leq r \left(\frac{2}{r} \right)^s = \frac{2^s}{r^{s-1}}.$$

Since each quadratic form q_k has a common variable with at most r other forms, we have

(3.6.2) For every k there are at most r indices $j \neq k$ such that $Q_k Q_j \neq 0$ or $Q_j Q_k \neq 0$.

Let $\omega \in \mathbb{C}$ be a complex number satisfying

$$(3.6.3) \quad |\omega| \leq \gamma = \frac{1}{4}e^{-\frac{1}{2}} \approx 0.1516326649.$$

Given Q_1, \dots, Q_k and ω , we construct a weighted graph $G = (V, E; w)$ as in Corollary 3.4. Our goal is to prove that $\text{Ind}_G(w) \neq 0$, for which we use Lemma 3.5.

We say that the *level* of a vertex $v = (k_1, \dots, k_s)$ is s for $s = 1, \dots, m$. Thus for the weight of v , we have

$$w(v) = \frac{\omega^s}{2^s} \text{trace}(Q_{k_1} \cdots Q_{k_s}).$$

Combining (3.6.1) and (3.6.3), we conclude that for a vertex of level s , we have

$$(3.6.4) \quad |w(v)| \leq \frac{1}{s2^{s+1}r^{s-1}}e^{-\frac{1}{2}s}.$$

We observe that there are at most sqr^{q-1} vertices u of level q with $w(u) \neq 0$ that are neighbors of a given vertex v (for $q = s$, we count v as its own neighbor). Indeed, there are at most s ways to choose a common index k , after which there are at most q positions to place k in u . By (3.6.2), we conclude that there are at most sqr^{q-1} vertices $u \in \mathcal{N}_v$ of level q with $w(u) \neq 0$. Choosing $\rho(v) = s/2$ for a vertex of level s and using (3.6.4), we conclude that for a vertex v of level s , we have

$$\begin{aligned} \sum_{u \in \mathcal{N}_v} |w(u)|e^{\rho(u)} &\leq \sum_{q=1}^m \left(\frac{1}{q2^{q+1}r^{q-1}}e^{-\frac{1}{2}q} \right) (sqr^{q-1}) e^{\frac{1}{2}q} \\ &= s \sum_{q=1}^m \frac{1}{2^{q+1}} < \frac{s}{2} = \rho(v), \end{aligned}$$

and the proof follows by Corollary 3.4. and Lemma 3.5. □

4. APPROXIMATING THE INTEGRAL

The interpolation method is based on the following simple observation.

(4.1) Lemma. *Let $p : \mathbb{C} \rightarrow \mathbb{C}$ be a polynomial,*

$$p(z) = \sum_{s=0}^m c_s z^s,$$

and $\beta > 1$ be a real number such that

$$p(z) \neq 0 \quad \text{provided} \quad |z| < \beta.$$

Let us choose a branch of $f(z) = \ln p(z)$ for $|z| < \beta$ and let

$$T_k(z) = f(0) + \sum_{s=1}^k \frac{f^{(s)}(0)}{s!} z^s$$

be the Taylor polynomial of degree k of f computed at $z = 0$. Then

$$|f(1) - T_k(1)| \leq \frac{m}{(k+1)\beta^k(\beta-1)}.$$

Moreover, the values of $f^{(s)}(0)$ for $s = 1, \dots, k$ can be computed from the coefficients c_s for $s = 0, \dots, k$ in time polynomial in n and m .

Proof. See, for example, Section 2.2 of [Ba16]. □

As follows from Lemma 4.1, if $\beta > 1$ is fixed in advance, to estimate the value of $f(1)$ within additive error $0 < \epsilon < 1$ (in which case we say that we estimate the value of $p(1) = e^{f(1)}$ within relative error ϵ), it suffices to compute the coefficients c_s with $s = O(\ln n - \ln \epsilon)$, where the implied constant in the “ O ” notation depends only on β . A similar result holds if $p(z) \neq 0$ in an arbitrary, fixed in advance, connected open set $U \subset \mathbb{C}$ such that $\{0, 1\} \subset U$, see Section 2.2 of [Ba16] (in Lemma 4.1, the neighborhood U is the disc of radius β).

(4.2) Computing the integrals. Let us fix a constant

$$0 < \gamma' < \gamma,$$

where γ is the constant of Theorem 1.3 (so one can choose $\gamma' = 0.15$). Let $q_1, \dots, q_m : \mathbb{R}^n \rightarrow \mathbb{R}$ be quadratic forms, defined by their matrices Q_1, \dots, Q_m as in (1.1.1), such that each form depends on not more than r variables among x_1, \dots, x_n and each form has common variables with not more than r other forms. Suppose that the bound (1.3.1) holds. We define a univariate polynomial $p : \mathbb{C} \rightarrow \mathbb{C}$ by

$$p(z) = \frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} (1 + zq_1(x)) \cdots (1 + zq_m(x)) e^{-\|x\|^2/2} dx.$$

Hence $\deg p \leq m$ and by Theorem 1.3 we have

$$p(z) \neq 0 \quad \text{provided} \quad |z| < \beta \quad \text{where} \quad \beta = \frac{\gamma}{\gamma'} > 1.$$

In view of Lemma 4.1, to approximate

$$(4.2.1) \quad p(1) = \frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} (1 + q_1(x)) \cdots (1 + q_m(x)) e^{-\|x\|^2/2} dx$$

within relative error $0 < \epsilon < 1$, it suffices to compute $p(0) = 1$ and $p^{(s)}(0)$ for $s = O(\ln m - \ln \epsilon)$, where the implied constant in the “ O ” notation is absolute. From Corollary 3.4, we have

$$p^{(s)}(0) = s! \sum_{\substack{(k_{11}, \dots, k_{1s_1}), \dots, (k_{j1}, \dots, k_{js_j}) \\ s_1 + \dots + s_j = s}} \frac{1}{2s_1} \cdots \frac{1}{2s_j} \\ \times \text{trace}(Q_{k_{11}} \cdots Q_{1s_1}) \cdots \text{trace}(Q_{k_{j1}} \cdots Q_{k_{j_s}}),$$

where the sum is taken over all unordered collections of pairwise disjoint ordered tuples $(k_{11}, \dots, k_{1s_1}), \dots, (k_{j1}, \dots, k_{js_j})$ of distinct indices k_{ij} from the set $\{1, \dots, m\}$, with the total number s of chosen indices. A crude upper bound for the number of such collections is $(2m)^s$: writing all the indices k_{ij} as a row, we have at most $2m$ choices for each index k_{ij} , including the choice on whether the index remains in the current tuple or starts a new one. Given that $s = O(\ln m - \ln \epsilon)$ and that computing the traces of the products of $n \times n$ matrices can be done in $(ns)^{O(1)}$ time, we obtain an algorithm approximating the integral in quasi-polynomial $n^{O(1)}m^{O(\ln m - \ln \epsilon)}$ time.

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