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VOLUME OF SUB-LEVEL SETS OF HOMOGENEOUS POLYNOMIALS

JEAN B. LASSERRE

Abstract. Consider the sub-level set \( K := \{ x : g(x) \leq 1 \} \) of a nonnegative homogeneous polynomial \( g \). We show that its Lebesgue volume \( \text{vol}(K) \) can be approximated as closely as desired by solving a sequence of generalized eigenvalue problems with respect to a pair of Hankel matrices of increasing size, whose entries are obtained in closed form. The methodology also extends to compact sets of the form \( \{ x : a \leq g(x) \leq b \} \) for non-homogeneous polynomials with degree \( d \ll n \). It reduces the volume computation in \( \mathbb{R}^n \) to a “volume” computation in \( \mathbb{R}^d \) (where \( d = \deg(g) \)) for a certain pushforward measure. Another extension to computing volumes of finite intersections of such sub-level sets is also briefly described.

MSC: 65K05 68U05 65D18 65D30 65F15 68W25 68W30 90C22

1. Introduction

Let \( g \in \mathbb{R}[x]_t \) be a nonnegative homogeneous polynomial of degree \( t \) (hence \( t \) is even) with associated sub-level set

\[
K := \{ x \in \mathbb{R}^n : g(x) \leq 1 \}.
\]

In this paper we describe an efficient numerical scheme to approximate its Lebesgue volume \( \text{vol}(K) \) (when finite) as closely as desired.

Motivation. In addition of being an interesting mathematical problem on its own, computing \( \text{vol}(K) \) has also a practical interest outside computational geometry. For instance it has a direct link with computing the integral

\[
\int_{\mathbb{R}^n} \exp(-g(x)) \, dx,
\]

called an integral discriminant in Dolotin and Morozov [3] and Morozov and Shakirov [11]. Indeed as proved in [11]:

\[
\text{vol}(K) = \frac{1}{\Gamma(1 + \frac{n+t}{2})} \int_{\mathbb{R}^n} \exp(-g(x)) \, dx,
\]

and to quote [11], “averaging with exponential weights is an important operation in statistical and quantum physics”. However, and again quoting [11], “despite simply looking, (1.2) remains terra incognita”. However, for special cases of homogeneous polynomials, the authors in [11] have been able to obtain a closed form expression for (1.2) (hence equivalently for \( \text{vol}(K) \)) in terms of algebraic invariants of \( g \).

Various consequences of formula (1.2) have been described and exploited in Lasserre [6]. For instance, \( \text{vol}(K) \) is a convex function in the coefficients of the...
matrices are easy to obtain in closed-form and the Hankel matrix $B$ of a monotone sequence of generalized eigenvalue problems with respect to a pair $(K, A)$. More precisely, $\text{vol}(K)$ is the smallest generalized eigenvalue of the pair $(A_d, B_d)$. The integral discriminant (1.2) is the limit of a monotone sequence of generalized eigenvalue problems with respect to a pair $(A_d, B_d)$ of given real Hankel matrices of size $d + 1$. All entries of both Hankel matrices are easy to obtain in closed-form and the Hankel matrix $B_d$ depends only on the degree of $g$. Therefore, in principle the integral (1.2) can be approximated efficiently and as closely as desired by (linear algebra) eigenvalue routines. To the best of our knowledge this result is quite new and in addition, even if we do not provide a closed form expression of (1.2), its new characterization as a limit or upper bound or a lower bound on $\text{vol}(K)$ is slow because of a Gibbs phenomenon.

Our main result is that

\begin{equation}
\text{vol}(K) = \lim_{d \to \infty} \lambda_{\min}(A_d, B_d),
\end{equation}

where $\lambda_{\min}(A_d, B_d)$ is the smallest generalized eigenvalue of the pair $(A_d, B_d)$. The integral (1.3) is a hard problem.

Methodology. Computing (and even approximating) the Lebesgue volume of a convex body is hard (let alone non-convex bodies). Often the only possibility is to use (non deterministic) Monte Carlo type methods which provide an estimate with statistical guarantees; that is, generate a sample of $N$ points according to the uniform distribution on $[-1, 1]^n$ and then the ratio $\rho_N := \text{(number of points in $K$)}/N$ provides such an estimate. However $\rho_N$ is a random variable and is neither an upper bound or a lower bound on $\text{vol}(K)$. For a discussion on volume computation the interested reader is referred to [4] and the many references therein.

However for basic semi-algebraic sets $K \subset [-1, 1]^n$, Henrion et al. [4] have provided a general methodology to approximate $\text{vol}(K)$. It consists in solving a hierarchy $(Q_d)_{d \in \mathbb{N}}$ of semidefinite programs\(^1\) of increasing size, whose associated sequence of optimal values $(\rho_d)_{d \in \mathbb{N}}$ is monotone non increasing and converges to $\text{vol}(K)$. An optimal solution of $Q_d$ is a vector $y \in \mathbb{R}^{s(2d)}$ (with $s(d) = \binom{n+d}{d}$) whose each coordinate $y_{s, \alpha} \in \mathbb{N}_d^n$, approximates the $\alpha$-moment of $\lambda_K$, the restriction to $K$ of the Lebesgue measure $\lambda$ on $\mathbb{R}^n$ (and therefore $y_0$ approximates $\text{vol}(K)$ from above).

An optimal solution of the dual semidefinite program $Q^*_d$ provides the coefficients $(p_n)_{n \in \mathbb{N}_d}$ of a polynomial $p \in \mathbb{R}[x]_{2d}$ which approximates on $[-1, 1]$ and from above, the (indicator) function $x \mapsto 1_K(x) = 1$ if $x \in K$ and 0 otherwise. In general the convergence $\rho_d \to \text{vol}(K)$ is slow because of a Gibbons' phenomenon\(^2\).

\(^1\)A semidefinite program (SDP) is a conic convex optimization problem with a remarkable modeling power. It can be solved efficiently (in time polynomial in its input size) up to arbitrary precision fixed in advance; see e.g. Anjos and Lasserre [1]

\(^2\)The Gibbs' phenomenon appears at a jump discontinuity when one approximates a piecewise $C^1$ function with a continuous function, e.g. by its Fourier series.
when one approximates the indicator function $1_K$ by continuous functions. In [4] the authors have proposed a “trick” which accelerate drastically the convergence but at the price of loosing the monotone convergence $\rho_d \downarrow \operatorname{vol}(K)$. Another acceleration technique was provided in [8] which still preserves monotone convergence. It uses the fact that moments of $\lambda_K$ satisfy linear equality constraints that follows from Stokes’ theorem.

Recently, Jasour et al. [5] have considered volume computation in the context of risk estimation in uncertain environments. They have provided an elegant “trick” which reduces computing the $n$-dimensional volume $\operatorname{vol}(K)$ to computing $\phi([0,1])$ for a certain pushforward measure $\phi$ on the real line, whose moments are known. (With $K$ as in (1.1) the pushforward measure $\phi$ is with respect to the mapping $g$.) This results in solving the hierarchy of semidefinite programs proposed in [4], but now for measures on the real line as opposed to measures on $\mathbb{R}^n$. Solving the corresponding hierarchy of dual semidefinite programs amounts to approximate the indicator of an interval on the real line by polynomials of increasing degree, and whose coefficients minimize a linear criterion.

On the one hand, it yields drastic computational savings as passing from $\mathbb{R}^n$ to $\mathbb{R}$ is indeed a big and impressive progress. But on the other hand the (monotone) convergence remains slow as one cannot one cannot apply the acceleration technique based on Stokes’ theorem proposed e.g. in [8] because the density of $\phi$ is not known explicitly. In the examples provided in §3.2 for comparison, we can observe this typical (very) slow convergence. However as the problem is now one-dimensional one may then solve many more steps of the resulting hierarchy of semidefinite programs provided that one works with a nice basis of polynomials (e.g. Chebyshev polynomials) to avoid numerical problems as much as possible. Interestingly, pushforward measures were also used in Magron et al. [10] to compute the Lebesgue volume of $f(K)$ for a polynomial mapping $f : \mathbb{R}^n \to \mathbb{R}^m$, but in this case one has to compute moments of the measure in $\mathbb{R}^n$ whose pushforward measure is the Lebesgue measure on $f(K)$, and the resulting computation is still very expensive and limited to modest dimensions.

**Contribution.** We provide a simple numerical scheme to approximate $\operatorname{vol}(K)$ with $K$ as in (1.1) and when $g$ is positive and homogeneous. To do so we are inspired by the trick of using the pushforward measure in Jasour et al. [5]. The novelty here is that by taking into account the specific nature (homogeneity) of $g$ in (1.1) we are able to drastically simplify computations. Indeed, the hierarchy of semidefinite programs defined in [5] can be replaced (and significantly improved) with computing a sequence of scalars $(\tau_d)_{d \in \mathbb{N}}$ where $\tau_d$ is nothing less than the generalized minimum eigenvalue of two known Hankel matrices of size $d$, whose entries are obtained exactly in closed-form with no numerical error. Therefore there is no optimization involved anymore. Moreover, if one uses the basis of orthonormal polynomials w.r.t. the pushforward measure, then $\tau_d$ is now the minimum eigenvalue of a single real symmetric matrix of size $d$.

At last but not least, the philosophy underlying the methodology also extends to arbitrary compact sets of the form $\{x : a \leq g(x) \leq b\} \subset \mathbb{R}^n$ where the polynomial $g$ is not necessarily homogenous and positive. This can be potentially interesting when $\deg(g) \ll n$ because we reduce the initial Lebesgue volume computation in $\mathbb{R}^n$ to a $\mu$-volume computation in $\mathbb{R}^d$ (where $d = \deg(g)$) for a certain pushforward measure $\mu$ on $\mathbb{R}^d$ whose sequence of moments is easily obtained in closed
form. Moreover, and crucial for the approach, we are still able to include additional constraints based on Stokes’ theorem, which significantly accelerates the otherwise typically slow convergence. In [5] the problem would be reduced to a ν-volume computation only in ℝ but with no possibility to include additional Stokes’ constraints to accelerate the slow convergence.

2. Notation and definitions

**Notation.** Let ℝ[x] denote the ring of polynomials in the variables x = (x₁, . . . , xₙ) and ℝ[x]ₜ ⊂ ℝ[x] denote the vector space of polynomials of degree at most t, hence of dimension s(d) = \( \binom{n+t}{n} \). Let Σ[x] ⊂ ℝ[x] denote the space of polynomials the are sums-of-squares (in short SOS polynomials) and let Σ[x]₂ᵈ ⊂ ℝ[x]₂ᵈ denote the space of SOS polynomials of degree at most 2d. With α ∈ ℕⁿ and x ∈ ℝⁿ, the notation x^α stands for \( x₁^α₁ \cdots xₙ^αₙ \). Also for every α ∈ ℕⁿ, let \( |α| := \sum_i α_i \) and \( ℕ^d_0 := \{ α ∈ ℕⁿ : |α| ≤ d \} \).

The support of a Borel measure μ on ℝⁿ is the smallest closed set Ω such that μ(Ω \ ∪) = 0. Denote by B(X) the Borel σ-field associated with a topological space X, and M(X) the space of finite Borel measures on X.

Given two real symmetric matrices A, C ∈ ℝⁿ×ⁿ denote by \( \lambda_{min}(A, C) \) the smallest generalized eigenvalue with respect to the pair (A, C), that is, the largest scalar θ such that \( Ax = θCx \) for some vector x ∈ ℝⁿ. When C is the identity matrix then \( \lambda_{min}(A, C) \) is just the smallest eigenvalue of A. Computing \( \lambda_{min}(A, C) \) can be done via a pure and efficient linear algebra routine. The notation \( A \succeq 0 \) (resp. \( A \succ 0 \)) stands for A is positive semidefinite (resp. positive definite).

**Moment matrix.** Given a real sequence \( φ = (φ_α)_{α ∈ ℕⁿ} \), let \( M_d(φ) \) denote the multivariate (Hankel-type) moment matrix defined by \( M_d(φ)(α, β) = φ_{α+β} \) for all \( α, β ∈ ℕ^d_0 \). For instance, in the univariate case \( n = 1 \), with \( d = 2 \), \( M_2 \) is the Hankel matrix

\[
M_2(φ) = \begin{bmatrix}
φ₀ & φ₁ & φ₂ \\
φ₁ & φ₂ & φ₃ \\
φ₂ & φ₃ & φ₄
\end{bmatrix}.
\]

If \( φ = (φ_j)_{j ∈ ℕ} \) is the moment sequence of a Borel measure φ on ℝ then \( M_d(φ) \succeq 0 \) for all \( d = 0, 1, . . . \). Conversely, if \( M_d(φ) \succeq 0 \) for all \( d ∈ ℕ \), then φ is the moment sequence of some finite Borel measure φ on ℝ. The converse result is not true anymore in the multivariate case.

Let φ, ν be two finite Borel measures on ℝ. The notation φ ≤ ν stands for φ(B) ≤ ν(B) for all \( B ∈ B(ℝ) \).

**Lemma 2.1.** Let φ, ν be two finite Borel measures on ℝ with all moments \( φ = (φ_j)_{j ∈ ℕ} \) and \( ν = (ν_j)_{j ∈ ℕ} \) finite. Then φ ≤ ν if and only if

\[
M_d(φ) ≤ M_d(ν), \ ∀ d = 0, 1, . . .
\]

**Proof.** Only if part: φ ≤ ν implies that ν - φ with associated sequence \( ν - φ = (ν_j - φ_j)_{j ∈ ℕ} \) is a finite Borel measure on ℝ, and therefore:

\[
M_d(ν) - M_d(φ) = M_d(ν - φ) ≥ 0, \ ∀ d ∈ ℕ,
\]

i.e., \( M_d(ν) ≥ M_d(φ) \) for all \( d ∈ ℕ \).

If part: If \( M_d(φ) ≤ M_d(ν) \) for all \( d ∈ ℕ \) then the sequence \( ν - φ = (ν_j - φ_j)_{j ∈ ℕ} \) satisfies \( M_d(ν - φ) ≥ 0 \) for all \( d ∈ ℕ \). Therefore, the moment sequence \( ν - φ = \)
Let \( A \) be the Lebesgue measure on \( \mathbb{R} \), and let \( M_d(p\phi) \) denote the real symmetric matrix defined by:

\[
M_d(p\phi)(\alpha, \beta) = \sum_{\gamma} p_\gamma \phi_{\alpha+\beta+\gamma}, \quad \alpha, \beta \in \mathbb{N}_d^2.
\]

For instance, with \( n = 1 \), \( d = 2 \) and \( x \mapsto p(x) = x(1-x) \):

\[
M_2(p\phi) = \begin{bmatrix}
\phi_1 - \phi_2 & \phi_2 - \phi_3 & \phi_3 - \phi_4 \\
\phi_2 - \phi_3 & \phi_3 - \phi_4 & \phi_4 - \phi_5 \\
\phi_3 - \phi_4 & \phi_4 - \phi_5 & \phi_5 - \phi_6
\end{bmatrix},
\]

also a Hankel matrix.

**Lemma 2.2.** Let \( x \mapsto p(x) = x(1-x) \).

(i) If a real (finite) sequence \( \phi = (\phi_j)_{j \leq 2d} \) satisfies \( M_d(\phi) \succeq 0 \) and \( M_{d-1}(p\phi) \succeq 0 \), then there is a measure \( \mu \) on \([0,1]\) whose moments \( \mu = (\mu_j)_{j \leq 2d} \) match \( \phi \).

(ii) If a real (infinite) sequence \( \phi = (\phi_j)_{j \in \mathbb{N}} \) satisfies \( M_d(\phi) \succeq 0 \) and \( M_d(p\phi) \succeq 0 \) for all \( d \), then there is a measure \( \mu \) on \([0,1]\) whose moments \( \mu = (\mu_j)_{j \in \mathbb{N}} \) match \( \phi \).

See for instance Lasserre [9] and the many references therein.

**Pushforward measure.** Let \( K \subset \mathbb{R}^n \) be a Borel set and \( \lambda \) a probability measure on \( K \). Given a measurable mapping \( f : K \to \mathbb{R}^p \), the pushforward measure of \( \lambda \) on \( \mathbb{R}^p \) w.r.t. \( f \) is denoted by \#\( \lambda \) and satisfies:

\[
\#\lambda(B) := \lambda(f^{-1}(B)), \quad \forall B \in \mathcal{B}(\mathbb{R}^p).
\]

In particular, its moments are given by

\[
\#\lambda_\alpha := \int_{\mathbb{R}^p} z^\alpha \#\lambda(dz) = \int_K f(x)^\alpha \lambda(dx), \quad \alpha \in \mathbb{N}^p.
\]

**A version of Stokes’ theorem.** Let \( \Omega \subset \mathbb{R}^n \) be an open subset with boundary \( \partial\Omega \) and let \( x \mapsto X(x) \) be a given vector field. Then under suitable smoothness assumptions,

\[
\int_\Omega \text{Div}(X)f(x) \, dx + \int_\Omega \langle X, \nabla f(x) \rangle \, dx = \int_{\partial\Omega} \langle \vec{n}_x, X \rangle f(x) \, d\sigma,
\]

where \( \vec{n}_x \) is the outward pointing normal to \( \Omega \) at \( x \in \partial\Omega \), and \( \sigma \) is the \((n-1)\)-dimensional Hausdorff measure on the boundary \( \partial\Omega \); see e.g. Taylor [12, Proposition 3.2, p. 128].

3. Main result

Let \( B := [-1,1]^n \) and \( K \subset B \) be as in (1.1) with \( \partial K \subset \{ x : g(x) = 1 \} \). Let \( \lambda \) be the Lebesgue measure on \( B \) normalized to a probability measure so that \( \text{vol}(K) = 2^n \lambda(K) \). Let \( g \) in (1.1) be a nonnegative and homogeneous polynomial of degree \( t \). That is, \( g(\lambda x) = \lambda^t g(x) \) for all \( \lambda \in \mathbb{R}, x \in \mathbb{R}^n \). Denote by:

\[
5_g := \max\{ g(x) : x \in B \}; \quad 0_g := \min\{ g(x) : x \in B \}.
\]

Notice that \( 0_g = 0 \) because \( 0 \in B \) and \( g \) is nonnegative with \( g(0) = 0 \), and therefore \( g(B) = [0, 5_g] \). We next follow an elegant idea of Jasour et al. [5], adapted to the
present context. It reduces the computation of $\text{vol}(K)$ (in $\mathbb{R}^n$) to a certain volume computation in $\mathbb{R}$, by using a particular pushforward measure of the Lebesgue measure $\lambda$ on $B$.

Let $\#\lambda$ be the pushforward on the positive half line of $\lambda$, by the polynomial mapping $g : B \to [0, \overline{b}_g]$. From (3.1), the support of $\#\lambda$ is the interval $I := [0, \overline{b}_g] \subset \mathbb{R}$. Then in view of (2.1):

$$\#\lambda_k := \int_I z^k \#\lambda(dz) = \int_B g(x)^k \lambda(dx), \quad k = 0, 1, \ldots$$

All scalars $(\#\lambda_k)_{k \in \mathbb{N}}$ can be obtained in closed form as $g$ is a polynomial and $\lambda$ is the (normalized) Lebesgue measure on $B$. Namely, writing the expansion

$$x \mapsto g(x)^k = \sum_{\alpha \in \mathbb{N}_d^k} g_{\alpha} x^\alpha,$$

for some coefficients $(g_{\alpha})$, one obtains:

$$\#\lambda_k = 2^{-n} \sum_{\alpha \in \mathbb{N}_d^k} g_{\alpha} \left( \prod_{i=1}^{n} \frac{(1 - (-1)^{\alpha_i+1})}{\alpha_i + 1} \right), \quad k = 0, 1, \ldots$$

Next observe that $2^{-n}\text{vol}(K) = \#\lambda(g(K))$ and note that $g(K) = [0,1]$. Therefore following the recipe introduced in Henrion et al. [4], and with $S := [0,1] \subset [0, \overline{b}_g]$:

$$\#\lambda(S) = \max_{\phi \in \mathcal{M}(S)} \{ \phi(S) : \phi \leq \#\lambda \}.$$  

Denote by $\phi^*$ the measure on the real line which is the restriction to $S \subset I$ of the pushforward measure $\#\lambda$. That is,

$$\phi^*(B) := \#\lambda(B \cap S), \quad \forall B \in \mathcal{B}(\mathbb{R}).$$

Then $\phi^*$ is the unique optimal solution of (3.4) and therefore $\phi^*(S) = \#\lambda(S)$; see e.g. Henrion et al. [4]. Then to approximate $\phi^*(S)$ from above, one possibility is to solve the hierarchy of semidefinite relaxations:

$$\rho_d = \max_{\phi} \{ \phi_0 : 0 \preceq M_d(\phi) \preceq M_d(\#\lambda) ; \ M_{d-1}(x(1-x) \phi) \succeq 0 \}$$

where $\phi = (\phi_{ij})_{i,j \leq 2d}$, $M_d(\#\lambda)$ is the (Hankel) moment matrix (with moments up to order $2d$) associated with the pushforward measure $\#\lambda$, and $M_d(\phi)$ (resp. $M_{d-1}(x(1-x) \phi)$) is the Hankel moment (resp. localizing) matrix (with moments up to order $2d$) associated with the sequence $\phi$ and the polynomial $x \mapsto p(x) = x(1-x)$; see §2. Indeed (3.6) is a relaxation of (3.4) and the sequence $(\rho_d)_{d \in \mathbb{N}}$ is monotone non increasing and converges to $\phi^*(S) = \#\lambda(S)$ from above; see e.g. [4].

The dual of (3.6) is the semidefinite program

$$\rho^*_d = \max_{p \in \mathbb{R}[x]_{2d}} \{ \int p d\#\lambda : \ p - 1 = \sigma_0 + \sigma_2 x(1-x) \ , \ p, \sigma_0 \in \Sigma[x]_d ; \ \sigma_1 \in \Sigma[x]_{d-1} \},$$

and if $K$ has nonempty interior then $\rho^*_d = \rho_d$.

This is the approach advocated by Jasour et al. [5] and indeed this reduction of the initial (Lebesgue) volume computation in $\mathbb{R}^n$

$$\text{vol}(K) = \max_{\phi \in \mathcal{M}(K)} \{ \phi(K) : \phi \leq \lambda \}$$
to instead compute \#\(\lambda([0, 1])\) (in \(\mathbb{R}\)) by solving (3.4) is quite interesting as it yields drastic computational savings; in fact, solving the multivariate analogues for (3.8) of the univariate semidefinite relaxations (3.6) for (3.4), becomes rapidly impossible even for moderate \(d\), except for problems of modest dimension (say e.g. \(n \leq 4\)).

However it is important to notice that in general the convergence \(\rho_d \downarrow \#\lambda(S)\) is very slow and numerical problems are expected for large values of \(d\). To partially remedy this problem the authors of [5] suggest to express moment and localizing matrices in (3.6) in the Chebyshev basis (as opposed to the standard monomial basis). This allows to solve a larger number of relaxations but it does not change the typical slow convergence. The trick based on Stokes' theorem used in [8] cannot be used here because the dominating (or reference) measure \#\(\lambda\) in (3.4) is not the Lebesgue measure \(\lambda\) anymore (as in (3.8)). On the other hand the trick to accelerate convergence used in [4] can still be used, that is, in (3.6) one now maximizes \(L_\phi(x(1-x)) = \phi_1 - \phi_2\) instead of \(\phi_0\). If \(\phi^d = (\phi^d_j)_{j \leq 2d}\) is an optimal solution of (3.6) then \(\phi^d \rightarrow \#\lambda(S)\) as \(d\) increases but one loses the monotone convergence from above.

In the sequel we show that in the particular case where \(g\) is positive and homogeneous then one can avoid solving the hierarchy (3.6) and instead solve a hierarchy of simple generalized eigenvalue problems with no optimization involved and with a much faster convergence.

### 3.1. Exploiting homogeneity.

**A crucial observation.** Recall that \(2^{-n}\text{vol}(\mathbf{K}) = \#\lambda(S)\). So let \(\phi^*\) be as in (3.5), and let \(\phi^* = (\phi^*_j)_{j \in \mathbb{N}}\) be its associated sequence of moments. Consider the vector field \(\mathbf{x} \mapsto \mathbf{X}(\mathbf{x}) := \mathbf{x}\). Then \(\text{Div}(\mathbf{X}) = n\). In addition, as \(g\) is homogeneous of degree \(t\) then by Euler's identity, \(\langle \mathbf{x}, \nabla g(\mathbf{x}) \rangle = tg(\mathbf{x})\) for all \(\mathbf{x} \in \mathbb{R}^n\). Recall that \(\mathbf{K} \subset \text{int}(\mathbf{B})\) and therefore \(g(\mathbf{x}) = 1\) for every \(\mathbf{x} \in \partial\mathbf{K}\). Next, for every \(j \in \mathbb{N}\), as \(g(\mathbf{x})^j = 1\) on \(\partial\mathbf{K}\) for all \(j \in \mathbb{N}\), Stokes' Theorem (2.2) yields:

\[
0 = \int_{\partial\mathbf{K}} \langle \nabla_x, \mathbf{x} \rangle (1 - g(\mathbf{x})^j) d\sigma \quad \text{[as } g(\mathbf{x}) = 1 \text{ on } \partial\mathbf{K}] \\
= n \int_{\mathbf{K}} (1 - g(\mathbf{x})^j) \lambda(d\mathbf{x}) + \int_{\mathbf{K}} \langle \mathbf{x}, \nabla(1 - g(\mathbf{x})^j) \rangle \lambda(d\mathbf{x}) \quad \text{[by Stokes]} \\
= n \lambda(\mathbf{K}) - (n + jt) \int_{\mathbf{K}} g(\mathbf{x})^j \lambda(d\mathbf{x}) \\
= n \#\lambda(S) - (n + jt) \int_{\partial(\mathbf{K})} z^j \#\lambda(dz) = n \phi^*_0 - (n + jt) \phi^*_j,
\]

so that we have proved:

**Lemma 3.1.** Let \(\phi^*\) be the Borel measure on \(\mathbb{R}\) which is the restriction to \(S = [0, 1]\) of the pushforward measure \(\#\lambda\) on \(I\). Then its moments \(\phi^* = (\phi^*_j)_{j \in \mathbb{N}}\) satisfy:

\[
\phi^*_0 = 2^{-n}\text{vol}(\mathbf{K}); \quad \phi^*_j := \frac{n}{n + j t} \phi^*_0, \quad j = 1, 2, \ldots
\]

Remarkably, Lemma 3.1 relates all moments of \(\phi^*\) to its mass \(\phi^*_0 = 2^{-n}\text{vol}(\mathbf{K})\) in very simple manner! However it now remains to compute \(\phi^*_0\).
Computing $\phi_0^*$. Define $M_d^*$ to be the Hankel (moment) matrix with entries:

$$M_d^*(k, \ell) := \frac{n}{n + (k + \ell - 2)t}, \quad k, \ell = 1, 2, \ldots,$$

so that $\phi_0^* M_d^* = M_d(\phi^*)$ for all $d \in \mathbb{N}$, where $M_d(\phi^*)$ is the Hankel moment matrix associated with the sequence $\phi^*$ (equivalently with the measure $\phi^*$).

Similarly, define $M_{d,x(1-x)}^*$ to be the Hankel matrix with entries:

$$M_{d,x(1-x)}^*(k, \ell) := \frac{n}{n + (k + \ell - 1)t} - \frac{n}{n + (k + \ell - 2)t}, \quad k, \ell = 1, 2, \ldots,$$

so that $\phi_0^* M_{d,x(1-x)}^* = M_d(x(1-x) \phi^*)$ is the localizing matrix associated with $\phi^*$ and the polynomial $x \mapsto x(1-x)$, for all $d \in \mathbb{N}$. As $\phi^*$ is supported on $[0, 1]$ then $\phi_0^* M_{d,x(1-x)}^* \succeq 0$ for all $d$, which in turn implies

$$M_{d,x(1-x)}^* \succeq 0, \quad \forall d \in \mathbb{N},$$

because $\phi_0^* > 0$.

**Theorem 3.2.** For each $d \in \mathbb{N}$, let $M_d^*$ be as in (3.10) and let $M_d(\#\lambda)$ be the Hankel moment matrix associated with $\#\lambda$ (hence with sequence of moments as in (3.3)). Then:

$$\phi_d^* = \lim_{d \to \infty} \lambda_{\text{min}}(M_d(\#\lambda), M_d^*),$$

i.e., $\phi_d^*$ is the limit of a sequence of minimum generalized eigenvalues associated with the pair $(M_d(\#\lambda), M_d^*))$, $d \in \mathbb{N}$.

**Proof.** For every $d \in \mathbb{N}$, let

$$\tau_d := \lambda_{\text{min}}(M_d(\#\lambda), M_d^*) = \max \{ \tau : \tau M_d^* \preceq M_d(\#\lambda) \},$$

which implies $\tau_d M_d^* \preceq M_d(\#\lambda)$. In addition, $\tau_d M_{d,x(1-x)}^* \succeq 0$ follows from (3.12). On the other hand, as $\phi^* \preceq \#\lambda$, we also have $\phi_0^* M_d^* = M_d(\phi^*) \preceq M_d(\#\lambda)$ for all $d \in \mathbb{N}$. Hence, $\phi_0^* \preceq \tau_d$ for all $d \in \mathbb{N}$, and the sequence $(\tau_d)_{d \in \mathbb{N}}$ is monotone non-increasing.

We next show that $\tau_d \downarrow \phi_0^*$ as $d \to \infty$. As $\tau_d \succeq \phi_0^*$ for all $d$, $\lim_{d \to \infty} \tau_d = \tau^* \succeq \phi_0^*$. Consider the sequence $\mu = (\mu_j)_{j \in \mathbb{N}}$ defined by:

$$\mu_j = \tau^* \frac{n}{n + jt}, \quad j \in \mathbb{N}.$$

Then from $\tau_d M_d^* \succeq 0$ for all $d$, and the convergence $\tau_d \to \tau^*$, we obtain

$$\tau^* M_d^* = M_d(\mu) \succeq 0, \quad \forall d \in \mathbb{N}.$$

For similar reasons

$$\tau^* M_{d,x(1-x)}^* = M_d(x(1-x) \mu) \succeq 0, \quad \forall d \in \mathbb{N}.$$

By Lemma 2.2, $\mu$ is the moment sequence of a measure $\mu$ supported on $[0, 1]$ with mass $\mu_0 = \mu([0, 1]) = \tau^*$, and by construction we also have $\mu \preceq \#\lambda$. Therefore $\mu \in M(S)$ is a feasible solution of (3.4) which implies $\mu([0, 1]) \leq \phi_0^*$. But on the other hand,

$$\phi_0^* \preceq \tau^* = \mu([0, 1]) \leq \phi_0^*,$$

which yields the desired result $\tau^* = \phi_0^*$. \hfill $\Box$

Therefore to approximate $\text{vol}(K)$ from above, one proceeds as follows. Start with $d = 1$ and then
• Compute all moments of \( \# \lambda \) up to order \( 2d \) by (3.3).
• Compute \( \tau_d := \lambda_{\min}(M_d(\# \lambda), M_d^*) \)
• set \( d = d + 1 \) and repeat.

This produces the required monotone sequence of upper bounds \( (\tau_d)_{d \in \mathbb{N}} \) on \( \phi_0^* \), which converges to \( \phi_0^* = 2^{-n} \mathrm{vol}(K) \) as \( d \) increases. Finally, the following result shows that \( \tau_d \leq \rho_d \).

**Proposition 3.3.** For each \( d \in \mathbb{N} \), let \( \rho_d \) (resp. \( \tau_d \)) be as in (3.6) (resp. (3.14)). Then \( \rho_d \geq \tau_d \).

**Proof.** Consider the sequence \( \mu = (\mu_j)_{j \leq 2d} \) defined by:

\[
\mu_j = \tau_d \frac{n}{n + j}, \quad j \leq 2d.
\]

Then from (3.14), \( \tau_d M_d^* = M_d(\mu) \) and therefore, \( 0 \preceq M_d(\mu) \preceq M_d(\# \lambda) \). Similarly \( \tau_d M_d-1_\lambda(1-x) = M_d-1_\lambda(x(1-x)) \mu \geq 0 \). In other words, the sequence \( \mu \) is a feasible solution of (3.6), which implies \( \mu_0 (= \tau_d) \leq \rho_d \). \( \square \)

Hence the above eigenvalue procedure (with no optimization involved) provides a monotone sequence of upper bounds on \( \phi_0^* \) that are better than the sequence of upper bounds \( (\rho_d)_{d \in \mathbb{N}} \) obtained by solving the hierarchy of semidefinite relaxations (3.6). Notice also that the matrix \( M_d^* \) depends only on the degree of \( g \) and not on \( g \) itself.

In fact there is a simple interpretation of this improvement. In Problem (3.4) and in its associated semidefinite relaxations (3.6), one may include the additional constraints

\[
(3.15) \quad \phi_j = n \phi_0 / (n + j), \quad \forall j \leq 2d,
\]

coming from Stokes’ theorem applied to \( \phi^* \); see Lemma 3.1. Indeed we are allowed to do that because \( \phi^* \) (which is the unique optimal solution of (3.4)) satisfies these additional constraints. If it it does not change the optimal value of (3.4) it changes that of (3.6) as it makes the corresponding relaxation stronger and therefore \( \tau_d \leq \rho_d \) for all \( d \).

**Remark 3.4.** (i) If one uses the basis of orthonormal polynomials with respect to the pushforward measure \( \# \lambda \) then the new moment matrix \( M_d(\# \lambda) \) (expressed in this basis) is the identity matrix and therefore \( \tau_d^{-1} \) is the maximum eigenvalue of the corresponding matrix \( M_d^* \) (also expressed in that basis). This basis of orthonormal polynomials can be obtained from the decomposition \( M_d(\# \lambda) = \mathbf{D} \mathbf{D}^T \) for triangular matrices \( \mathbf{D} \) and \( \mathbf{D}^T \). For more details on multivariate orthogonal polynomials, the interested reader is referred to Denkhl and Xu [2] and the many references therein.

(ii) Alternatively one may also use an orthonormal basis associated with the sequence of moments \( (n/(n + j))_{j \in \mathbb{N}} \). In this case the new moment matrix \( M_d^* \) expressed in this basis is the identity matrix and \( \tau_d \) is now the minimum eigenvalue of the new moment \( M_d(\# \lambda) \) expressed in this basis. Notice that the orthonormal basis does not depend on \( g \) (only on its degree).

(iii) Finally, another possibility is to simply use the basis of Chebyshev polynomials. Then computing \( \tau_d \) (still a generalized eigenvalue problem) involves matrices with much better numerical conditioning.
3.2. Some numerical examples. To show how this approximation of \( \text{vol}(K) \) from above by a sequence of eigenvalue problems of increasing size is much more efficient than solving the hierarchy of semidefinite programs (3.6) as suggested in [5], we have considered a favorable case for (3.6). We chose \( K \) to be the Euclidean unit ball \( \{ x : \| x \| \leq 1 \} \) with Lebesgue volume \( \pi^{n/2}/\Gamma(1 + n/2) \) and the ball \( B \) that contains \( K \) is the smallest one, i.e., \( B = [-1,1]^n \). Indeed, the smaller is the ball \( B \), the better are the approximation by \( \rho_d \) in (3.6).

We first describe the first two steps to appreciate the simplicity of the approach.

Let \( n = 2 \) and \( g = \| x \|^2 = x_1^2 + x_2^2 \), and \( B = [-1,1]^2 \), so that \( \text{vol}(K) = \pi \). Then:

\[
M_1^* = \begin{bmatrix}
1 & 1/2 \\
1/2 & 1/3
\end{bmatrix}; \quad M_1(#\lambda) = \begin{bmatrix}
1 & 2/3 \\
2/3 & 28/45
\end{bmatrix}
\]

This yields \( 4 \cdot \tau_1 \approx 3.48 \) which is already a good upper bound on \( \pi \) whereas \( 4 \cdot \rho_1 = 4. 

\[
M_2^* = \begin{bmatrix}
1 & 1/2 & 1/3 \\
1/2 & 1/3 & 1/4 \\
1/3 & 1/4 & 1/5
\end{bmatrix}; \quad M_2(#\lambda) = \begin{bmatrix}
1 & 2/3 & 28/45 \\
2/3 & 28/45 & 24/35 \\
\end{bmatrix}
\]

This yields \( 4 \cdot \tau_2 \approx 3.1440 \) and \( 4 \cdot \rho_2 = 3.8928 \). Hence \( 4 \tau_2 \) already provides a very good upper bound on \( \pi \) with only moments of order 4. To appreciate the difference in speed of convergence between \( \rho_d \) and \( \tau_d \), Table 1 displays both values \( \tau_d \) and \( \rho_d \) in the case of \( n = 4 \) variables and \( d = 1, \ldots, 5 \). While the convergence \( \tau_d \to 4.9348 \) is quite fast with a relative error of 0.03\% at step \( d = 5 \), the convergence \( \rho_d \to 4.9348 \) is extremely slow as \( \rho_5 \approx 8.499 \) only; see Figure 1.

![Figure 1. n = 4; Comparing \( \tau_d \) (red below) with \( \rho_d \) (blue above)](image)

We next provide results for the same problem but now in larger dimensions \( n = 5, 8, 9, 10 \) in Table 3, Table 4, Table 5, and Table 6 respectively. From inspection
we can observe a fast and regular decrease in the value $2^n \tau_d$ as $d$ increases, and similarly for the relative error.

For $n = 10$ and $d = 8$, we have encountered numerical problems because the Hankel matrix $M_s(\# \lambda)$ is ill-conditioned and then one should use another basis of polynomials in which to express the matrices $M_d$ and $M_s(\# \lambda)$; see Remark 3.4.

**Influence of the size of the ball $B$.** If one increases the size of the ball $B = [-r, r]^n$ that contains $K$ then one expects a slower convergence and this is why it is recommended to take for $B$ the smallest ball that contains $K$. An appropriate choice is the box $\prod_{i=1}^n [-u_i, u_i]$ where $u_i$ (resp. $v_i$) is a lower bound (resp. upper bound) as close as possible to $\min \{ x_i : x \in K \}$ (resp. $\max \{ x_i : x \in K \}$), which can be computed by the first step of the Moment-SOS hierarchy described in [9]. From results displayed in Table 7 with $r = 1$ and $r = 1.3$, one observes that even though the convergence is a bit slower it is still quite good. The initial value $\tau_1$ is significantly higher but then $\tau_d$ (with $r = 1.3$) still decreases very fast; see Figure 2.

4. Extensions

In this section we discuss two extensions of the above methodology, when:

- $K := \{ x : a \leq g(x) \leq b \} \subset (-1, 1)^n$ and $g$ is not homogeneous anymore.
- $K$ is now $\{ x : g_j(x) \geq 0, j = 1, \ldots, m \} \subset (-1, 1)^n$ and each $g_j$ is homogeneous (with one of them being nonnegative).
In the second extension, again following Jasour et al. [5], one considers the pushforward of the Lebesgue measure by the polynomial mapping $g : \mathbf{B} \to \mathbb{R}^m$ which maps $\mathbf{x} \in \mathbf{B}$ to the vector $(g_j(\mathbf{x}))_{j=1}^{\deg(g)} \in \mathbb{R}^m$. Then the initial Lebesgue volume computation in $\mathbb{R}^n$ is reduced to an equivalent “measure” computation problem of the form (3.8) but now in $\mathbb{R}^m$ (instead of $\mathbb{R}$) and for which one may apply the hierarchy of semidefinite programs described in [4]. But as we did in §3, we can exploit again the homogeneity of the $g_j$’s to strengthen the semidefinite relaxations defined in [5], by introducing additional linear constraint coming from an appropriate application of Stokes’ theorem. The only difference with the univariate case treated in §3 is that the problem is not an eigenvalue problem anymore.

The first extension to the non homogeneous case is perhaps more interesting. We now write $g$ as a sum of homogeneous polynomials of increasing degree $1, 2, \ldots, \deg(g)$, and consider again a pushforward of the Lebesgue measure $\lambda$ by the polynomial mapping $g : \mathbf{B} \to \mathbb{R}^{\deg(g)}$, $\mathbf{x} \mapsto (g_1(\mathbf{x}), \ldots, g_{\deg(g)}(\mathbf{x}))$. 

In the second extension, again following Jasour et al. [5], one considers the pushforward of the Lebesgue measure by the polynomial mapping $g : \mathbf{B} \to \mathbb{R}^m$ which maps $\mathbf{x} \in \mathbf{B}$ to the vector $(g_j(\mathbf{x}))_{j=1}^{\deg(g)} \in \mathbb{R}^m$. Then the initial Lebesgue volume computation in $\mathbb{R}^n$ is reduced to an equivalent “measure” computation problem of the form (3.8) but now in $\mathbb{R}^m$ (instead of $\mathbb{R}$) and for which one may apply the hierarchy of semidefinite programs described in [4]. But as we did in §3, we can exploit again the homogeneity of the $g_j$’s to strengthen the semidefinite relaxations defined in [5], by introducing additional linear constraint coming from an appropriate application of Stokes’ theorem. The only difference with the univariate case treated in §3 is that the problem is not an eigenvalue problem anymore.

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\begin{table}[h]
\centering
\begin{tabular}{|c|c|c|c|c|c|c|}
\hline
$d$ & $d = 2$ & $d = 3$ & $d = 4$ & $d = 5$ & $d = 6$ & $d = 7$ \\
\hline
$2^n \tau_d$ & 32.432 & 12.657 & 6.662 & 4.375 & 3.379 & 2.921 \\
$\frac{(\tau_d - \rho^*)}{\rho^*}$ & 1171% & 396.3% & 161% & 71.6% & 32.5% & 14.54% \\
\hline
\end{tabular}
\caption{Table 6. $n = 10$, $\rho^* = 2.55$; $\tau_d$ and relative error}
\end{table}

\begin{table}[h]
\centering
\begin{tabular}{|c|c|c|c|c|c|c|c|}
\hline
$d$ & $d = 1$ & $d = 2$ & $d = 3$ & $d = 4$ & $d = 5$ & $d = 6$ & $d = 7$ & $d = 8$ \\
\hline
$r = 1.3$: $(2r)^n \tau_d$ & 20.345 & 11.744 & 7.622 & 6.149 & 5.585 & 5.373 & 5.299 & 5.275 \\
$r = 1$: $(2r)^n \tau_d$ & 10.289 & 6.524 & 5.575 & 5.334 & 5.278 & 5.266 & 5.264 & 5.2639 \\
\hline
\end{tabular}
\caption{Table 7. $n = 5$; $\text{vol}(\mathbf{K}) = 5.2638$; Influence of the size of $\mathbf{B} = [-r, r]^n$ with $r = 1$ and $r = 1.3$}
\end{table}

Figure 2. $n = 8$; Comparing $\tau_d$ with $r = 1.3$ (red above) and $r = 1$ (blue below)
4.1. The non-homogeneous case. Let $B = [-1, 1]^n$, and suppose that $K \subset \mathbb{R}^n$ is now described by:

\begin{equation}
K := \{x : a \leq g(x) \leq b\},
\end{equation}

for some $a, b \in \mathbb{R}$, where $g \in \mathbb{R}[x]$, and $K \subset (-1, 1)^n$, possibly after scaling. With no loss of generality we may and will assume that $g(0) = 0$ and write

\[ x \mapsto g(x) = \sum_{k=1}^{t} g_{k}(x), \quad x \in \mathbb{R}^n, \]

where for each $1 \leq k \leq t$, $g_{k}$ is a homogeneous polynomial of degree $k$.

We next see how to adapt the previous methodology of §3 to this more general case in a relatively simple manner. To simplify the exposition and alleviate notation, we will describe the quadratic case $t = 2$ but it will become obvious to understand how to proceed for $t > 2$. So with $t = 2$, $g = g_{1} + g_{2}$ with $g_{1}$ (resp. $g_{2}$) homogeneous of degree 1 (resp. 2).

Consider the pushforward measure $\#\lambda$ on $\mathbb{R}^2$ of $\lambda$ on $B$, by the polynomial mapping:

\[ g : B \to \mathbb{R}^2, \quad x \mapsto g(x) = \begin{bmatrix} g_{1}(x) \\ g_{2}(x) \end{bmatrix}, \quad x \in B. \]

Let $\Theta := g(B) \subset \mathbb{R}^2$ be the support of the pushforward measure $\#\lambda$, and observe that for each $i, j \in \mathbb{N}$:

\begin{equation}
\#\lambda_{ij} := \int_{\Theta} z_{1}^{i} z_{2}^{j} d\#\lambda(z) = \int_{B} g_{1}(x)^{i} g_{2}(x)^{j} d\lambda(x),
\end{equation}

can be obtained in closed form. Letting

\[ S := g(K) = \{z \in \Theta : a \leq z_{1} + z_{2} \leq b\}, \]

we obtain $\text{vol}(K) = \#\lambda(S)$. Next, recall that (see (3.4) in §3):

\begin{equation}
\text{vol}(K) = \#\lambda(S) = \max_{\phi \in \mathcal{M}(S)} \{ \phi(S) : \phi \preceq \#\lambda \}
\end{equation}

and $\phi^{*}$ is the unique optimal solution of (4.3).

Let $z \mapsto \tilde{h}(z) := (b - z_{1} - z_{2})(z_{1} + z_{2} - a)$. The semidefinite relaxations associated with (4.3) read:

\begin{equation}
\rho_d = \max_{\phi} \{ \phi_0 : 0 \preceq M_d(\phi) \preceq M_d(\#\lambda) ; \quad M_d(\tilde{h} \phi) \succeq 0 \},
\end{equation}

where the maximization is over finite bivariate sequences $\phi = (\phi_{ij})_{i+j \leq 2d}$.

Following the same philosophy as in §3, we are going to use some additional information on the optimal solution $\phi^{*}$ of (4.3) to strengthen the semidefinite relaxations (4.4). To do so we again use Stokes’ theorem.

**Stokes.** Recall that $K \subset (-1, 1)^n$ and therefore, $\partial K \subset \{x \in B : h(x) = 0\}$ where $x \mapsto h(x) := (b - g_{1}(x) - g_{2}(x))(g_{2}(x) + g_{2}(x) - a)$. Therefore by Stokes’ theorem,

\[ 0 = n \int_{K} g_{1}(x)^{i} g_{2}(x)^{j} h(x) d\lambda(x) \]

\[ + \int_{K} \langle x, \nabla(g_{1}(x)^{i} g_{2}(x)^{j} h(x)) d\lambda(x), \quad \forall i, j \in \mathbb{N}. \]
Developing and using homogeneity of \( g_1, g_2 \), one obtains:

\[
0 = n \int_S \left[ (n + i + 2j) z_1^i z_2^j (b - z_1 - z_2)(z_1 + z_2 - a) \right. \\
+ (z_1 + 2z_2) z_1^i z_2^j (a + b - 2z_1 - 2z_2) \bigg]. \, d\lambda(z).
\]

Equivalently, introducing the polynomial \( q_{ij} \in \mathbb{R}[z] \),

\[
(4.6) \quad z \mapsto q_{ij}(z) := (n + i + 2j) z_1^i z_2^j (b - z_1 - z_2)(z_1 + z_2 - a) \\
+ (z_1 + 2z_2) z_1^i z_2^j (a + b - 2z_1 - 2z_2).
\]

for every \( i, j \in \mathbb{N} \), one obtains:

\[
(4.7) \quad \int_S q_{ij}(z) \, d\phi^*(z) = 0, \quad \forall i, j \in \mathbb{N}.
\]

Notice that (4.7) is a linear relation between moments of \( \phi^* \), the optimal solution of (4.3). That is, let \( \phi^* = (\phi^*_{ij})_{i,j \in \mathbb{N}} \) be the sequence of moments of \( \phi^* \) on \( S \), and let \( L_{\phi^*} : \mathbb{R}[z] \to \mathbb{R} \) be the Riesz functional

\[
q = \sum_{i,j} q_{ij} z_1^i z_2^j \mapsto L_{\phi^*}(q) := \sum_{i,j} q_{ij} \phi^*_{ij}, \quad q \in \mathbb{R}[z].
\]

Then (4.7) reads

\[
(4.8) \quad L_{\phi^*}(q_{ij}) = 0, \quad i, j \in \mathbb{N},
\]

So we can strengthen the relaxations (4.4) by adding the additional (Stokes) moments constraints (4.8), that is, one solves the semidefinite programs:

\[
(4.9) \quad \tau_d = \max_{\phi} \{ \phi_0 : 0 \preceq M_d(\phi) \preceq M_d(\#\lambda); \quad M_d(\tilde{h} \phi) \succeq 0; \quad L_{\phi^*}(q_{ij}) = 0, \ \text{for all} \ (i, j) \ \text{s.t. deg}(q_{ij}) \leq 2d \},
\]

which is clearly a strengthening of (4.4).

**Proposition 4.1.** Let \( \rho_d \) (resp. \( \tau_d \)) be as in (4.4) (resp. (4.9)), \( d \in \mathbb{N} \). Then:

\[
(4.10) \quad \#\lambda(S) \leq \tau_d \leq \rho_d \quad \text{for all} \ d, \ \text{and} \ \tau_d \downarrow \#\lambda(S) \ \text{as} \ d \ \text{increases}.
\]

**Proof.** That \( \tau_d \leq \rho_d \) for all \( d \), is straightforward and similarly for the monotonicity of the sequence \( (\tau_d)_{d \in \mathbb{N}} \). Next, as \( \phi^* \) is the optimal solution of (4.3), its sequence of moments \( \phi^* = (\phi^*_{ij}) \) is feasible for (4.9), with associated value \( \phi^*_{ij} = \#\lambda(S) = \text{vol}(K) \). Hence \( \tau_d \geq \#\lambda(S) \). Then the convergence \( \tau_d \downarrow \#\lambda(S) \) follows from \( \rho_d \downarrow \#\lambda(S) \).

The difference with the homogeneous case treated in §3 is that now computing \( \tau_d \) requires solving the semidefinite program (4.9) whereas in §3 computing \( \tau_d \) reduces to solving a generalized eigenvalue problem, hence with no optimization involved. However, notice that instead of solving the (costly) \( n \)-variate semidefinite relaxations associated with (3.8) in \( \mathbb{R}^n \), we now solve similar semidefinite relaxations but for a bivariate problem on the plane. In addition the (convergence) acceleration technique based on Stokes’s theorem can also be implemented (see (4.9)).
4.2. Multi-homogeneous constraints. Another extension is when $K = \{ \mathbf{x} : g_j(\mathbf{x}) \leq 1, j = 1, \ldots, m \} \subset (-1, 1)^n$ for a family $(g_j)_{j=1}^m$ of homogeneous polynomials, not necessarily of same degree, say $\deg(g_j) = t_j$, and at least one of them is positive on $(\mathbb{R} \setminus \{0\})^n$. In this case one may proceed again as suggested in Jasour et al. [5]. Now $\# \lambda$ is the pushforward on $\mathbb{R}^m$ of $\lambda$ on $B$, by the mapping:
\[
g : B \to \mathbb{R}^m, \quad g(\mathbf{x}) = \begin{bmatrix} g_1(\mathbf{x}) \\ \vdots \\ g_m(\mathbf{x}) \end{bmatrix}.
\]
In particular $\# \lambda$ has its moments defined by:
\[
\# \lambda_\alpha = \int_B g_1(\mathbf{x})^{\alpha_1} \cdots g_m(\mathbf{x})^{\alpha_m} \lambda(d\mathbf{x}) = \int_{g(B)} z^\alpha \# \lambda(dz), \quad \forall \alpha \in \mathbb{N}^m.
\]
Again all moments $\# \lambda_\alpha$ can be computed in closed form, and again with $S = [0, 1]^m$ $2^{-n} \text{vol}(K) = \# \lambda(S)$. Let us describe how the generalization works for the case $m = 2$. Again denote by $\phi^*$ on $\mathbb{R}^2$ the restriction of $\# \lambda$ to $S$ and let $\phi^* = (\phi^*_i)_{i,j \in \mathbb{N}}$ with:
\[
\phi^*_{ij} := \int_S z_1^i z_2^j \phi^*(dz), \quad \forall i, j = 0, 1, \ldots.
\]
So the bivariate analogues of the semidefinite relaxations (3.6) read:
\[
\rho_d = \max_{\phi} \{ \phi_0 : 0 \preceq M_d(\phi) \preceq M_d(\# \lambda) \}
\]
\[
(4.11)
\]
where $\phi = (\phi_{ij})_{i+j \leq 2d}$, and $M_d(\phi)$ (resp. $M_{d-1}(x_j(1-x_j) \phi)$, $j = 1, 2$) is the moment (resp. localizing) matrix associated with $\phi$ (resp. with $\phi$ and $\mathbf{x} \mapsto x_j(1-x_j)$, $j = 1, 2$). Then $\rho_d \downarrow \# \lambda(S)$ as $d \to \infty$. Again the semidefinite relaxations (4.11) are a lot cheaper to solve than those associated with the $n$-variate problem (3.8). As we did for the univariate case we can improve the above convergence by adding additional constraints that must be satisfied at the optimal solution $\phi^*$ of (3.8). Again $\phi_0^* = 2^{-n} \text{vol}(K)$. Let $(i, j, k, \ell) \in \mathbb{N}^4$ with $k, \ell \geq 1$. Then with $X(\mathbf{x}) = \mathbf{x}$, Stokes’s Theorem yields
\[
0 = n \int_K g_1^i g_2^j (1-g_1)^k (1-g_2)^\ell \lambda(d\mathbf{x})
\]
\[
+ \int_K \langle \mathbf{x}, \nabla [g_1^i g_2^j (1-g_1)^k (1-g_2)^\ell] \rangle \lambda(d\mathbf{x})
\]
\[
= n \int_S z_1^i z_2^j (1-z_1)^k (1-z_2)^\ell \# \lambda(dz)
\]
\[
+ it_1 \int_S z_1^i z_2^j (1-z_1)^k (1-z_2)^\ell \# \lambda(dz)
\]
\[
+ jt_2 \int_S z_1^i z_2^j (1-z_1)^k (1-z_2)^\ell \# \lambda(dz)
\]
\[
- kt_1 \int_S z_1^{i+1} z_2^j (1-z_1)^{k-1} (1-z_2)^\ell \# \lambda(dz)
\]
\[
- \ell t_2 \int_S z_1^i z_2^{j+1} (1-z_1)^k (1-z_2)^{\ell-1} \# \lambda(dz)
\]
That is, for each \((i, j, k, \ell) \in \mathbb{N}^4\) with \(k, \ell \geq 1\) one obtains a linear constraint that links the moments of \(\phi^*\), that we denote by \(L_{\phi^*}(g_{ijk\ell})\) where \(g_{ijk\ell} \in \mathbb{R}[x]\) is the above polynomial under the integral sign. For instance,

\[
0 = L_{\phi^*}(q_{0011}) = n(\phi_0^* - \phi_{10}^* - \phi_{01}^* + \phi_{11}^*) - t_1(\phi_0^* - \phi_{01}^*) - t_2(\phi_0^* - \phi_{10}^*),
\]

\[
0 = L_{\phi^*}(q_{1111}) = (n + t_1 + t_2)(\phi_{11}^* + \phi_{22}^* - \phi_{12}^* - \phi_{21}^*) - t_1(\phi_{21}^* - \phi_{12}^*) - t_2(\phi_{12}^* - \phi_{22}^*),
\]

etc. So we can add these additional constraints to (4.11) and solve:

\[
\tau_d = \max_{\phi} \{ \phi_0 : 0 \leq M_d(\phi) \leq M_d(\#\lambda) \}
\]

(4.12)

Of course \(\tau_d \leq \rho_d\) for all \(d\) and therefore \(\tau_d \downarrow \#\lambda(S)\) as \(d\) increases. Again, the difference with the univariate case is that now computing \(\tau_d\) requires to solve the semidefinite program (4.12) instead of a generalized eigenvalue problem. However it is of same dimension as (4.11) and the convergence \(\tau_d \downarrow \#\lambda(S)\) is expected to be much faster than \(\rho_d \downarrow \#\lambda(S)\) as we have been able to include additional constraints based on Stokes’ theorem.

5. Conclusion

We have presented a new methodology to approximate (in principle as closely as desired) the Lebesgue volume of the sub-level set \(\{x : g(x) \leq 1\}\) of a positive homogeneous \(n\)-variate polynomial \(g\). Inspired by Jasour et al. [5], we formulate an equivalent “volume” computation \(\mu(I)\) of an interval \(I\) of the real line for a certain pushforward measure \(\mu\). The novelty with respect to [5] is that by using Stokes’ theorem and exploiting the homogeneity of \(g\), we are able to further reduce the problem to a hierarchy of generalized eigenvalue problems for Hankel matrices of increasing size, with no optimization involved. To the best of our knowledge, this characterization of Lebesgue volume as the limit of eigenvalue problems of increasing size is new. Moreover, the methodology also extends to sub-level sets of arbitrary polynomials. It then reduces the Lebesgue volume computation in \(\mathbb{R}^n\) to a “volume” computation in \(\mathbb{R}^d\) (where \(d\) is the degree of the initial polynomial) for a certain pushforward measure. An extension to the case several homogeneous constraints has been also presented with a similar pattern as in the extension to the single non-homogeneous constraint. Preliminary results on a simple case reveal a drastic improvement on the approximation scheme proposed in [5].

References


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