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A transfer theorem for multivariate Δ -analytic functions with a power-law singularity

Linxiao Chen

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Abstract

This paper presents a multivariate generalization of Flajolet and Odlyzko's transfer theorem. Similarly to the univariate version, the theorem assumes Δ -analyticity (defined coordinate-wise) of a function $A(z_1, \ldots, z_d)$ at a unique dominant singularity $(\rho_1, \ldots, \rho_d) \in (\mathbb{C}^*)^d$, and allows one to translate, on a term-by-term basis, an asymptotic expansion of $A(z_1, \ldots, z_d)$ around (ρ_1, \ldots, ρ_d) into a corresponding asymptotic expansion of its Taylor coefficients a_{n_1,\ldots,n_d} . We treat the case where the asymptotic expansion of $A(z_1, \ldots, z_d)$ contains only power-law type terms, and where the indices n_1, \ldots, n_d tend to infinity in some polynomially stretched diagonal limit. The resulting asymptotic expansion of a_{n_1,\ldots,n_d} is a sum of terms of the form

$$I(\lambda_1,\ldots,\lambda_d)\cdot n_0^{-\Theta}\cdot \rho_1^{-n_1}\cdots \rho_d^{-n_d}$$

where $(\lambda_1, \ldots, \lambda_d) \in (0, \infty)^d$ is the direction vector of the stretched diagonal limit for (n_1, \ldots, n_d) , the parameter n_0 tends to ∞ at similar speed as n_1, \ldots, n_d , while $\Theta \in \mathbb{R}$ and $I : (0, \infty)^d \to \mathbb{C}$ are determined by the asymptotic expansion of A.

1 Introduction

Univariate transfer theorems. In their pioneering work [5], Flajolet and Odlyzko developped a method, known as the transfer theorems, that allows one to compute a precise asymptotic expansion of a sequence a_n as $n \to \infty$, from an asymptotic expansion of its generating function $A(z) = \sum_{n\geq 0} a_n z^n$ around some singularity $\rho \in \mathbb{C}^* \equiv \mathbb{C} \setminus \{0\}$. We use the notation $[z^n]A(z)$ for the coefficient $a_n = \frac{1}{n!}A^{(n)}(0)$. One distinctive feature of the transfer theorem is that it applies to generating functions that are Δ -analytic, that is, analytic on a Δ -domain of the form $\rho \cdot \Delta_{\delta} \equiv \{\rho z \mid z \in \Delta_{\delta}\}$ for some $\delta > 0$, where

$$\Delta_{\delta} := \left\{ z \in \mathbb{C} : |z| < 1 + \delta, z \neq 1 \text{ and } \arg(1-z) \in \left(-\frac{\pi}{2} - \delta, \frac{\pi}{2} + \delta\right) \right\}.$$

$$\tag{1}$$

See Figure 1(a). (Simple extensions of the transfer theorem also apply to functions analytic in a finite intersection of Δ -domains of the form $(\rho_1 \cdot \Delta_\delta) \cap \cdots \cap (\rho_n \cdot \Delta_\delta)$, but we shall not discuss that further here.) Under the Δ -analyticity assumption, ρ is the unique dominant singularity of the function A(z). By the change of variable $z' = z/\rho$, we can bring this singularity to z = 1. In the following, we will assume without loss of generality that $\rho = 1$.

In its most general form, Flajolet and Odlyzko's transfer theorem applies to functions whose asymptotic expansion is composed of any regular varying functions taken from a large class of "standard scale functions", such as $f(z) = (1-z)^{\alpha} (\log(1-z))^{\beta}$ for $\alpha, \beta \in \mathbb{R}$. Here we focus on functions whose asymptotic expansion contains only power-law terms of the form $(1-z)^{\alpha}$. In this simple case, the transfer theorem can be formulated as follows.

Theorem A (Univariate transfer theorem). Let $\alpha \in \mathbb{R}$ and let A(z) be an analytic function on Δ_{δ} for some $\delta > 0$.

- 1. (Analytic terms have exponentially small contribution) If A(z) is also analytic at 1, then there exists $r \in (0,1)$ such that $[z^n]A(z) = O(r^n)$ when $n \to \infty$. (In this case, z = 1 is actually not the dominant singularity of A(z).)
- 2. (Coefficient asymptotics of power functions) When $A(z) = (1-z)^{\alpha}$ with $\alpha \in \mathbb{R}$, we have the following asymptotic expansion as $n \to \infty$,

$$[z^n](1-z)^{\alpha} = \frac{n^{-\alpha-1}}{\Gamma(-\alpha)} \sum_{k=0}^{\infty} \frac{e_k(\alpha)}{n^k} \equiv \frac{n^{-\alpha-1}}{\Gamma(-\alpha)} \left(1 + \frac{e_1(\alpha)}{n} + \frac{e_2(\alpha)}{n^2} + \cdots \right)$$
(2)

where $e_k(\alpha) \in \mathbb{Q}[\alpha]$ is a polynomial of degree 2k in α . More precisely, $e_k(\alpha) = \sum_{l=0}^k g_{k,l} \frac{\Gamma(-\alpha)}{\Gamma(-\alpha-k-l)}$, where $g_{k,l}$ is defined by the Taylor expansion

$$G(x,y) := \exp\left(-y - \left(\frac{y}{x} + 1\right)\log(1-x)\right) = \sum_{k=0}^{\infty} \left(\sum_{l=0}^{k} g_{k,l} y^{l}\right) x^{k}.$$
 (3)

- 3. (Big-O transfer and little-o transfer)
 - If $A(z) = O((1-z)^{\alpha})$ when $z \to 1$ in Δ_{δ} , then $[z^n]A(z) = O(n^{-\alpha-1})$ when $n \to \infty$.
 - If $A(z) = o((1-z)^{\alpha})$ when $z \to 1$ in Δ_{δ} , then $[z^n]A(z) = o(n^{-\alpha-1})$ when $n \to \infty$.

In practice, the transfer theorem is usually applied to functions A(z) which are linear combinations of the three cases above, as in the following statement.

Theorem A' (Univariate transfer theorem, integrated form). Let A(z) be an analytic function on Δ_{δ} for some $\delta > 0$ such that when $z \to 1$ in Δ_{δ} , we have

$$A(z) = A_{\rm reg}(z) + h_0(1-z)^{\alpha_0} + \dots + h_m(1-z)^{\alpha_m} + o\left((1-z)^{\alpha_m}\right) \,, \tag{4}$$

where $A_{reg}(z)$ is defined and analytic in a neighborhood of 1 in \mathbb{C} , and $h_j \in \mathbb{C}$, $\alpha_j \in \mathbb{R}$ for all $0 \leq j \leq m$. Then, the coefficients of A(z) has the following expansion when $n \to \infty$

$$[z^{n}]A(z) = h_{0} \cdot [z^{n}](1-z)^{\alpha_{0}} + \dots + h_{m} \cdot [z^{n}](1-z)^{\alpha_{m}} + o(n^{-\alpha_{m}-1})$$
(5)

where each term $[z^n](1-z)^{\alpha_j}$ has the asymptotic expansion given in Theorem A(2). The same result holds if the little-o estimates in both (4) and (5) are replaced by big-O estimates.

Remark.

- Theorem A and Theorem A' are not exactly equivalent, for the following reasons: When A(z) is analytic at 1, Theorem A asserts that [zⁿ]A(z) decays exponentially as n → ∞, while Theorem A' only implies a super-polynomial decay. On the other hand, one cannot prove Theorem A' by simply applying Theorem A to each term of the expansion (4), because the terms A_{reg}(z) and o((1 α)^{αm}) in (4) have a priori no analytic continuations on the whole domain Δ_δ.
- 2. For a sequence of functions $s_k(n)$, the asymptotic expansion $S(n) = \sum_{k=0}^{\infty} s_k(n)$ usually means that for all $m \ge 0$, we have $S(n) = s_0(n) + \cdots + s_{m-1}(n) + O(s_m(n))$ as $n \to \infty$. Theorem A(2) uses a slightly modified notion of asymptotic expansion: we choose implicitly the family $s_{\beta}(n) = n^{-\beta}$ $(\beta \in \mathbb{R})$ as the reference asymptotic scale, and view each asymptotic expansion as a weighted sum of the form $S(n) = \sum_{k=0}^{\infty} c_k \cdot n^{-\beta_k}$ with some $\beta_0 < \beta_1 < \cdots$. The difference is that now the prefactors c_k may vanish for some or even all $k \ge 0$, and the expansion should be read as: for all $m \ge 0$, we have $S(n) = c_0 n^{-\beta_0} + \cdots + c_{m-1} n^{-\beta_{m-1}} + O(n^{-\beta_m})$ as $n \to \infty$.
- 3. It is also possible to include $\alpha \in \mathbb{C} \setminus \mathbb{R}$ in Theorem A and theorem A'. The asymptotic expansion (2) would hold unmodified. But complex exponents α rarely appear in applications and they complicate the asymptotic scale $s_{\beta}(n) = n^{-\beta}$ (since \mathbb{C} is not totally ordered). We restrict ourselves to $\alpha \in \mathbb{R}$ for the sake of simplicity.

4. Theorem A(1)-(2) together imply that the prefactor $\frac{1}{\Gamma(-\alpha)}$ in (2) must vanish when $\alpha \in \mathbb{N}$. Indeed, $\alpha \mapsto \frac{1}{\Gamma(-\alpha)}$ is an entire function and we have $\frac{1}{\Gamma(-\alpha)} = 0$ if and only if $\alpha \in \mathbb{N}$.

The transfer theorems apply to a large class of naturally occuring generating functions. For example, it is well known that all D-finite functions analytic at 0 are linear combinations of Δ -analytic functions. Among them, the algebraic functions always have singularities of power-law type, to which Theorem A' applies. Compared to alternative methods such as Darboux's method or the Tauberian theorems, the transfer theorem has the advantage of giving a transparent correspondance between the asymptotic expansion of the generating function and that of its coefficients. See [6, Sec. 5] for a detailed discussion about this comparison.

In this paper, we present a generalization of the transfer theorem to the multivariate setting. Below, we start by recalling some basic notations about multivariate generating functions, and then define the regimes of coefficient asymptotics with which our transfer theorem will be concerned.

Multivariate generating functions. In many practical problems, the relevant information is captured naturally by a multidimensional infinite array $(a_n)_{n \in \mathbb{N}^d} \in \mathbb{C}^{\mathbb{N}^d}$. Such multidimensional arrays and the multivariate generating functions which encode them will be this paper's central objects. To make the formulas compact, we will use the following multi-index notations: For any formal or complex vectors $\boldsymbol{z} = (z_1, \ldots, z_d)$ and $\boldsymbol{\theta} = (\theta_1, \ldots, \theta_d)$, we denote

$$\boldsymbol{\theta}\boldsymbol{z} = (\theta_1 z_1, \dots, \theta_d z_d), \qquad \boldsymbol{\theta} \cdot \boldsymbol{z} = \theta_1 z_1 + \dots + \theta_d z_d, \qquad \boldsymbol{z}^{\boldsymbol{\theta}} = z_1^{\theta_1} \cdots z_d^{\theta_d}, \qquad \mathrm{d}\boldsymbol{z} = \mathrm{d} z_1 \cdots \mathrm{d} z_d.$$
(6)

And, for any scalar σ and integer vector $\boldsymbol{m} \in \mathbb{N}^d$, let

$$\sigma^{\boldsymbol{\theta}} = (\sigma^{\theta_1}, \dots, \sigma^{\theta_d}) \quad \text{and} \quad \boldsymbol{\partial}^{\boldsymbol{m}} = \partial_1^{m_1} \cdots \partial_d^{m_d} \equiv \frac{\partial^{m_1}}{\partial z_1^{m_1}} \cdots \frac{\partial^{m_d}}{\partial z_d^{m_d}}.$$
 (7)

With these notations, the multivariate generating function of the array $(a_n)_{n \in \mathbb{N}^d}$ can be written as

$$A(\boldsymbol{z}) = \sum_{\boldsymbol{n} \in \mathbb{N}^d} a_{\boldsymbol{n}} \boldsymbol{z}^{\boldsymbol{n}} \equiv \sum_{n_1=0}^{\infty} \cdots \sum_{n_d=0}^{\infty} a_{n_1,\dots,n_d} z_1^{n_1} \cdots z_d^{n_d}.$$
(8)

Similarly to the univariate case, we denote the coefficient a_n by $[\boldsymbol{z}^n]A(\boldsymbol{z})$. We assume that every generating function in this paper is absolutely convergent in an open neighborhood of $\mathbf{0} \equiv (0, \ldots, 0) \in \mathbb{C}^d$, so that it defines an analytic function there.

We refer to [7] for the general theory on power series and analytic functions in several variables. One particular fact that we will use without further mention is the uniqueness of analytic continuation: if A and B are analytic functions on an open connected domain $\Omega \subseteq \mathbb{C}^d$ and A(z) = B(z) on any open subset of Ω , then A(z) = B(z) for all $z \in \Omega$.

Stretched diagonal limits. One central problem of analytic combinatorics in several variables is to understand the asymptotics of the coefficients $[\boldsymbol{z}^n]A(\boldsymbol{z})$ when the components n_1, \ldots, n_d of the multi-index tend to ∞ simultaneously. In general, one needs to put some constraint on the relative speeds at which n_1, \ldots, n_d grow in order to get a useful asymptotic formula, and the interesting regimes are to a large extent dictated by the structure of the singularities of $A(\boldsymbol{z})$.

In this work, we will be interested in the *stretched diagonal limits*, where n_1, \ldots, n_d grow at polynomial speeds relative to each other. In other words, we require that for each $j = 2, \ldots, d$, there exists a constant $\theta_j > 0$ such that the ratio $n_j/n_1^{\theta_j}$ remains in some compact interval $\mathcal{I} \subset \mathbb{R}_{>0} \equiv (0, \infty)$ when $n_1 \to \infty$. A symmetrized definition of this limit regime goes as follows:

Definition (Stretched diagonal limit). We say that the multi-index $n \in \mathbb{N}^d$ tends to ∞ in the *stretched* diagonal regime if there exist an exponent vector $\boldsymbol{\theta} \in \mathbb{R}^d_{>0}$ and an auxiliary variable $n_0 > 0$, such that

$$\boldsymbol{n} = \boldsymbol{\lambda} n_0^{\boldsymbol{\theta}} \equiv \left(\lambda_1 n_0^{\theta_1}, \dots, \lambda_d n_0^{\theta_d} \right)$$
(9)

for some prefactor $\lambda = \lambda(n_0)$ that remains in a compact set $\mathcal{K} \subset \mathbb{R}^d_{>0}$ when $n_0 \to \infty$. In this case, we will also say that $n \to \infty$ in the θ -diagonal limit.

Remark. If we require in addition that $\lambda(n_0) \to \lambda_* \in \mathbb{R}^d_{>0}$ as $n_0 \to \infty$, then the point $n = \lambda n_0^{\theta}$ would tend to ∞ roughly along the curve $\mathcal{D}_{\lambda_*,\theta} = (\lambda_* t^{\theta}, t > 0)$. But we only require $\lambda(n_0)$ to stay in some compact set. Intuitively this means that n can jump between the curves $\mathcal{D}_{\lambda_*,\theta}$ for different values of λ_* . Consequently, the asymptotics that we write in the stretched diagonal limit should be understood as uniform with respect to λ_* on every compact subset of $\mathbb{R}^d_{>0}$.

Notice that for any $\tau > 0$, the $(\tau \theta)$ -diagonal limit is the same as the θ -diagonal limit: it suffices to replace the variable n_0 by n_0^{τ} or $n_0^{1/\tau}$ to pass from one to the other. The classical notion of *diagonal limit*, e.g. as defined in [9], corresponds to the $(1, \dots, 1)$ -diagonal limit in our terminology.

Apart from being the relevant limit regime for multivariate Δ -analytic function with power-law singularity (to be defined below), the stretch diagonal limits also arise naturally from the study of critical phenomena in probability and mathematical physics.

Preliminary definitions. Let us define some domains and function classes needed for stating the main theorems. Let $\mathbb{C}_* = \mathbb{C} \setminus \{0\}$ and $\mathbb{R}_{>0} = (0, \infty)$. For $\delta > 0$, we write $K_{\delta} = \{u \in \mathbb{C}_* : |\arg(u)| < \delta\}$ and denote $\Omega_{\delta} = K_{\pi/2+\delta}$. The multivariate versions of these cones are denoted $\mathbf{K}_{\delta} = K_{\delta}^d$ and $\Omega_{\delta} = \Omega_{\delta}^d$. For $\epsilon > 0$, $w \in \mathbb{C}$ and $w \in \mathbb{C}^d$, let $B_{\epsilon}(w) = \{z \in \mathbb{C} : |z - w| < \epsilon\}$ and $B_{\epsilon}(w) = B_{\epsilon}(w_1) \times \cdots \times B_{\epsilon}(w_d)$. Notice that Ω_{δ} is related to the Δ -domain Δ_{δ} by $\Omega_{\delta} = \{\mu(1-z) \mid \mu > 0 \text{ and } z \in \Delta_{\delta}\}$.

Definition (Multivariate Δ -analytic functions). Let $\boldsymbol{\rho} \in \mathbb{C}^d_*$. We say that a multivariate function $A(\boldsymbol{z})$ is Δ -analytic at $\boldsymbol{\rho}$ if it has an analytic continuation on the product domain $\boldsymbol{\rho} \Delta_{\delta} := (\rho_1 \Delta_{\delta}) \times \cdots \times (\rho_d \Delta_{\delta})$ for some $\delta > 0$, where Δ_{δ} is the univariate Δ -domain defined in (1).

Like in the univariate case, one can make the change of variable $\mathbf{z}' = \mathbf{z}/\boldsymbol{\rho} \equiv (z_1/\rho_1, \dots, z_d/\rho_d)$ to bring the point $\boldsymbol{\rho}$ to $\mathbf{1} \equiv (1, \dots, 1)$. In the following, we will focus without loss of generality on functions which are Δ -analytic at $\mathbf{1}$.

Definition (Demi-analytic functions). Let $\Omega = \Omega_1 \times \cdots \times \Omega_d$ be an open product domain in \mathbb{C}^d and let $\rho \in (\partial \Omega_1) \times \cdots \times (\partial \Omega_d)$. For $j \in \{1, \ldots, d\}$, denote $\Omega_{\hat{j}} = \Omega_1 \times \cdots \times \Omega_{j-1} \times \mathbb{C} \times \Omega_{j+1} \times \cdots \times \Omega_d$. We say that a function $A : \Omega \to \mathbb{C}$ is *demi-analytic at* $\rho \in \Omega$ if A is analytic on Ω and there exist $\epsilon > 0$ and a decomposition $A = A_1 + \cdots + A_d$ such that A_j is analytic in $B_{\epsilon}(\rho) \cap \Omega_{\hat{j}}$ for each $j \in \{1, \ldots, d\}$. If each term A_j in the above decomposition is analytic on $\Omega_{\hat{j}}$, then we say that A is *demi-entire*.

Definition (Generalized homogeneous functions). Let K be a cone in \mathbb{C}^d , i.e. a subset of \mathbb{C}^d such that $\{\sigma \boldsymbol{z} \mid \boldsymbol{z} \in \boldsymbol{K}\} = \boldsymbol{K}$ for all $\sigma > 0$. For $\theta_0 \in \mathbb{R}$ and $\boldsymbol{\theta} \in \mathbb{R}^d_{>0}$, we say that a function $H : \boldsymbol{K} \to \mathbb{C}$ is $(\theta_0, \boldsymbol{\theta})$ -homogeneous if

$$H(\sigma^{\theta} \boldsymbol{u}) \equiv H(\sigma^{\theta_1} u_1, \dots, \sigma^{\theta_d} u_d) = \sigma^{\theta_0} H(\boldsymbol{u})$$
(10)

for all $\boldsymbol{u} \in \boldsymbol{K}$ and $\sigma > 0$.

It is clear that a (θ_0, θ) -homogeneous function is also $(\tau \theta_0, \tau \theta)$ -homogeneous for all $\tau > 0$. The classical notion of homogeneous functions of degree D becomes $(D, \mathbf{1})$ -homogeneous in our terminology.

The three definitions above generalize respectively the notions of Δ -analytic functions, locally analytic functions (for the term A_{reg} in Theorem A'), and power functions (for the term $h \cdot (1-z)^{\alpha}$) used in the univariate transfer theorems. To state the multivariate transfer theorem, we need one more definition whose counterpart does not appear explicitly in the univariate setting:

Definition (Functions of polynomial type). For any $\delta > 0$, we say that a function F is of polynomial type (globally) on \mathbf{K}_{δ} if there exist C, M > 0 such that

$$\forall \boldsymbol{u} \in \boldsymbol{K}_{\delta}, \quad |F(\boldsymbol{u})| \leq C \cdot \left(|u_1|^{-M} + |u_1|^M + \dots + |u_d|^{-M} + |u_d|^M\right)$$
(11)

We say that F is of polynomial type locally at $\mathbf{0} \in \mathbf{K}_{\delta}$ if the above bound only holds in a neighborhood of $\mathbf{0} \in \mathbf{K}_{\delta}$. Equivalently, F is of polynomial type locally at $\mathbf{0} \in \mathbf{K}_{\delta}$ if there exist $\epsilon, C, M > 0$ such that

$$\forall \boldsymbol{u} \in \boldsymbol{K}_{\delta} \cap \boldsymbol{B}_{\boldsymbol{0},\epsilon}, \quad |F(\boldsymbol{u})| \leq C \cdot \left(|u_1|^{-M} + \dots + |u_d|^{-M}\right).$$
(12)

Similarly, we say that a function A is of polynomial type on Δ_{δ} if there exist C, M > 0 such that

$$\forall \boldsymbol{z} \in \boldsymbol{\Delta}_{\delta}, \quad |A(\boldsymbol{z})| \le C \cdot \left(|z_1 - 1|^{-M} + \dots + |z_d - 1|^{-M}\right)$$
(13)

(here we do not need the terms $|z_j - 1|^M$ on the right hand side because Δ_{δ} is bounded), and we say that A is of polynomial type locally at $\mathbf{1} \in \Delta_{\delta}$ if it satisfies the above bound in $\Delta_{\delta} \cap B_{\mathbf{1},\epsilon}$ for some $\epsilon > 0$. For $S = K_{\delta}$, Ω_{δ} or Δ_{δ} , we write $\mathcal{P}(S) = \{f : S \to \mathbb{C} \mid f \text{ is analytic and of polynomial type on } S\}$.

Up to decreasing δ , a continuous function on Δ_{δ} is always bounded by a constant on $(\Delta_{\delta} \setminus U)^d$ for any neighborhood U of $1 \in \mathbb{C}$. Thus, the condition that A is of polynomial type (globally) on Δ_{δ} is essentially an upper bound for $|A(\mathbf{z})|$ when $z_j \to 1$ in Δ_{δ} for some $j \in \{1, \ldots, d\}$. In contrast, the condition of being of polynomial type locally at $\mathbf{1} \in \Delta_{\delta}$ only gives a bound for $|A(\mathbf{z})|$ when $z_j \to 1$ for all $j \in \{1, \ldots, d\}$. When d = 1, the two conditions are essentially the same. They do not appear explicitly in the statement of the univariarte transfer theorem because a function having an asymptotic expansion of the form (4) is automatically of polynomial type on Δ_{δ} .

Main results. We are now ready to state the multivariate transfer theorem. Like in the univariate case, we formulate it in two ways, one discussing the building blocks of the coefficient asymptotics piece by piece, and the other showing how the theorem would be applied in practice.

Theorem 1 (Multivariate transfer theorem). Let $\theta_0 \in \mathbb{R}$, $\boldsymbol{\theta} \in \mathbb{R}^d_{>0}$ and $A \in \mathcal{P}(\boldsymbol{\Delta}_{\delta})$ for some $\delta > 0$. We consider the asymptotics of the coefficients $[\boldsymbol{z}^n]A(\boldsymbol{z})$ when $\boldsymbol{n} \to \infty$ in the $\boldsymbol{\theta}$ -diagonal regime.

- 1. (Demi-analytic terms have exponentially small contribution) If A is demi-analytic at $\mathbf{1} \in \mathbf{\Delta}_{\delta}$, then there exist $r \in (0, 1)$ and $\sigma > 0$ such that $[\mathbf{z}^n]A(\mathbf{z}) = O(r^{n_0^{\sigma}})$.
- 2. (Coefficient asymptotics of generalized homogeneous functions) When $A(\mathbf{z}) = H(\mathbf{1} - \mathbf{z})$ and H is $(\theta_0, \boldsymbol{\theta})$ -homogeneous $(\theta_0 \in \mathbb{R})$, we have the asymptotic expansion

$$[\boldsymbol{z}^{\boldsymbol{n}}]H(\boldsymbol{1}-\boldsymbol{z}) = \frac{1}{n_0^{\Theta}} \sum_{\boldsymbol{k} \in \mathbb{N}^d} \frac{D_{\boldsymbol{k}}I(\boldsymbol{\lambda})}{n_0^{\boldsymbol{k}\cdot\boldsymbol{\theta}}} \equiv \frac{1}{n_0^{\Theta}} \left(I(\boldsymbol{\lambda}) + \sum_{j=1}^d \frac{\partial_j I(\boldsymbol{\lambda}) + \frac{1}{2}\lambda_j \partial_j^2 I(\boldsymbol{\lambda})}{n_0^{\theta_j}} + \cdots \right).$$
(14)

where $\Theta = \theta_0 + \theta_1 + \dots + \theta_d$, $D_{\mathbf{k}} = \sum_{l \leq \mathbf{k}} g_{\mathbf{k},l} \cdot \lambda^l \partial^{\mathbf{k}+l}$ is a partial differential operator of order $2(k_1 + \dots + k_d)$, and $I : \mathbb{R}^d_{>0} \to \mathbb{C}$ is the inverse Laplace transform of H defined by

$$I(\boldsymbol{\lambda}) = \left(\frac{1}{2\pi i}\right)^d \int_{\boldsymbol{V}_{\delta'}} H(\boldsymbol{u}) e^{\boldsymbol{\lambda} \cdot \boldsymbol{u}} \mathrm{d}\boldsymbol{u} \,. \tag{15}$$

In the above formulas, we denote $g_{\mathbf{k},\mathbf{l}} = g_{k_1,l_1} \cdots g_{k_d,l_d}$, with the numbers $(g_{k,l})_{k,l\geq 0}$ defined by (3). The sum $\sum_{\mathbf{l}\leq \mathbf{k}}$ runs over $\{\mathbf{l}\in \mathbb{N}^d \mid l_j\leq k_j \text{ for all } 1\leq j\leq d\}$, while the integral is over $V_{\delta'} = V_{\delta'}^d$, where $V_{\delta'} \subset \Omega_{\delta}$ is any piecewise smooth curve which coincides with the rays $\partial\Omega_{\delta'}$ outside a bounded region for some $\delta' \in (0, \delta)$, e.g. as in Figure 1(b). (Recall that $\Omega_{\delta} = \{u \in \mathbb{C}_* : |\arg(u)| < \frac{\pi}{2} + \delta\}$.)

3. (Big-O transfer and little-o transfer)

If $A(\boldsymbol{z}) = O(\tilde{H}(\boldsymbol{1}-\boldsymbol{z}))$ for some $(\theta_0, \boldsymbol{\theta})$ -homogeneous \tilde{H} as $\boldsymbol{z} \to \boldsymbol{1}$ in $\boldsymbol{\Delta}_{\delta}$, then $[\boldsymbol{z}^n]A(\boldsymbol{z}) = O(n_0^{-\Theta})$. If $A(\boldsymbol{z}) = o(\tilde{H}(\boldsymbol{1}-\boldsymbol{z}))$ for some $(\theta_0, \boldsymbol{\theta})$ -homogeneous \tilde{H} as $\boldsymbol{z} \to \boldsymbol{1}$ in $\boldsymbol{\Delta}_{\delta}$, then $[\boldsymbol{z}^n]A(\boldsymbol{z}) = o(n_0^{-\Theta})$. **Theorem 1'** (Multivariate transfer theorem). Let $\boldsymbol{\theta} \in \mathbb{R}^d_{>0}$ and $A \in \mathcal{P}(\boldsymbol{\Delta}_{\delta})$ for some $\delta > 0$. Assume that when $\boldsymbol{z} \to \boldsymbol{1}$ in $\boldsymbol{\Delta}_{\delta}$, the function A has an asymptotic expansion of the form

$$A(z) = A_{reg}(z) + H_0(1-z) + \dots + H_{m-1}(1-z) + o(H_m(1-z)) , \qquad (16)$$

where A_{reg} is demi-analytic at $\mathbf{1} \in \mathbf{\Delta}_{\delta}$, and each H_k is $(\theta_0^{(k)}, \boldsymbol{\theta})$ -homogeneous and of polynomial type locally at $\mathbf{0} \in \mathbf{\Omega}_{\delta}$, with $\theta_0^{(0)} > \cdots > \theta_0^{(m-1)} \ge \theta_0^{(m)}$. Then as $\mathbf{n} \to \infty$ in the $\boldsymbol{\theta}$ -diagonal limit, we have

$$[\boldsymbol{z}^{\boldsymbol{n}}]A(\boldsymbol{z}) = [\boldsymbol{z}^{\boldsymbol{n}}]H_0(\boldsymbol{1} - \boldsymbol{z}) + \dots + [\boldsymbol{z}^{\boldsymbol{n}}]H_{m-1}(\boldsymbol{1} - \boldsymbol{z}) + o(n_0^{-\Theta_m})$$
(17)

where each $[\boldsymbol{z}^n]H_k(1-\boldsymbol{z})$ has the asymptotic expansion given in Theorem 1(2), $\Theta_m = \theta_0^{(m)} + \theta_1 + \dots + \theta_d$, and the little-o estimate is uniform with respect to $\boldsymbol{\lambda}$ on all compact subsets of $\mathbb{R}^d_{>0}$.

The same result holds if the little-o estimates in both (16) and (17) are replaced by big-O estimates.

Remark. The following remarks are analoguous to their univariate counterparts below Theorem A'.

- 1. Theorems 1 and 1' are not exactly equivalent, for the same reason as in the univariate case.
- 2. The asymptotic expansion (14) in Theorem 1(2) uses the same modified notion of asymptotic expansion as in Theorem A(2), with the asymptotic scale $s_{\beta}(n_0) = n_0^{-\beta}$ ($\beta \in \mathbb{R}$).
- 3. It is also possible to include $\theta_0 \in \mathbb{C} \setminus \mathbb{R}$ in Theorems 1 and 1'. The expansion (14) would still hold with a complex value for Θ . But we restrict ourselves to $\theta_0 \in \mathbb{R}$ for simplicity.
- 4. Theorem 1(1)–(2) imply that for a (θ_0, θ) -homogeneous function $H \in \mathcal{P}(\Omega_{\delta})$, if H is demi-analytic at $\mathbf{0} \in \Omega_{\delta}$, then its inverse Laplace transform $I(\boldsymbol{\lambda})$ vanishes for all $\boldsymbol{\lambda} \in \mathbb{R}^d_{>0}$. We will see in Corollary 3 below that the converse is also true.

Let us also make some remarks about the expansion (16) in Theorem 1'. These will be discussed in more details in Section 4.

- 5. The expansion (16) is in general not unique. The reason is that a demi-analytic function can also contain (θ_0, θ) -homogeneous components for any $\theta_0 \in \mathbb{R}$. In principle, one could fix $A_{reg} = 0$ in Theorem 1' without reducing significantly the class of functions A(z) covered by the theorem. But having the flexibility of choosing any demi-analytic function A_{reg} makes the theorem easier to apply.
- 6. Once A_{reg} is chosen, one can write $A(\boldsymbol{z}) = A_{\text{reg}}(\boldsymbol{z}) + A_{\text{sing}}(1-\boldsymbol{z})$. If the expansion (16) exists, then its terms H_k can be obtained as the coefficients in the *univariate* asymptotic expansion $A_{\text{sing}}(\varepsilon^{\boldsymbol{\theta}}\boldsymbol{u}) = H_0(\boldsymbol{u}) \cdot \varepsilon^{\theta_0^{(0)}} + \cdots + H_{m-1}(\boldsymbol{u}) \cdot \varepsilon^{\theta_0^{(m-1)}} + o(\varepsilon^{\theta_0^{(m)}})$ with respect to $\varepsilon \to 0^+$.
- 7. Although not assumed in Theorem 1', the functions $\boldsymbol{z} \mapsto H_k(\boldsymbol{1} \boldsymbol{z})$ $(0 \le k < m)$ are necessarily analytic on $\boldsymbol{\Delta}_{\delta}$. See Lemma 4 below. In particular, the coefficients $[\boldsymbol{z}^n]H_k(\boldsymbol{1} \boldsymbol{z})$ are well-defined.
- 8. Unlike the univariate case, the generalized homogenous function H_m used in the little-o estimate is in general different from the previous term H_{m-1} of the asymptotic expansion. The reason is that the class of (θ_0, θ) -homogeneous functions is one-dimensional (generated by $H(u) = u^{\theta_0/\theta}$) in the univariate case, but infinite-dimensional in the multivariate case.

In many applications, one is only interested in the dominant asymptotics of the coefficients. For this, we can simplify Theorems 1 and 1' to the following statement: for any generating function $A \in \mathcal{P}(\Delta_{\delta})$,

$$A(\boldsymbol{z}) = A_{\text{reg}}(\boldsymbol{z}) + H(\boldsymbol{1} - \boldsymbol{z}) + o(\tilde{H}(\boldsymbol{1} - \boldsymbol{z})) \qquad \Rightarrow \qquad [\boldsymbol{z}^{\boldsymbol{n}}]A(\boldsymbol{z}) \sim I(\boldsymbol{\lambda}) \cdot n_0^{-\Theta}$$
(18)

where A_{reg} is demi-analytic and of polynomial type locally at $\mathbf{1} \in \Delta_{\delta}$, H and H are (θ_0, θ) -homogeneous and of polynomial type locally at $\mathbf{0} \in \Delta_{\delta}$, and $\Theta \in \mathbb{R}$ and $I : \mathbb{R}_{>0} \to \mathbb{C}$ are defined as in Theorem 1(2). As in the theorems, the asymptotics of the coefficients is taken in the θ -diagonal regime.

The functional form of the prefactor $I(\lambda)$ is often of practical importance (see Section 4 for more discussions). In Theorem 1, $I(\lambda)$ was expressed as an integral transform of H(u). The next theorem

will provide some basic properties of this transform and its inverse. We define the *inverse Laplace* transform (also called *Borel transform*) of a function H by

$$\mathcal{B}[H](\boldsymbol{\lambda}) = \left(\frac{1}{2\pi i}\right)^d \int_{\boldsymbol{V}_{\boldsymbol{\delta}'}} e^{\boldsymbol{\lambda} \cdot \boldsymbol{u}} H(\boldsymbol{u}) \,\mathrm{d}\boldsymbol{u}$$
(19)

where $\delta' \in (0, \delta)$ and $V_{\delta'}$ is a contour of the form specified below Equation (15). On the other hand, for any given $\boldsymbol{c} \in \mathbb{R}^d_{>0}$, we define the Laplace transform (truncated at \boldsymbol{c}) of a function I by

$$\mathcal{L}_{\boldsymbol{c}}[I](\boldsymbol{u}) = \int_{\boldsymbol{c}}^{\infty} e^{-\boldsymbol{\lambda} \cdot \boldsymbol{u}} I(\boldsymbol{\lambda}) \, \mathrm{d}\boldsymbol{\lambda} \equiv \int_{[c_1,\infty) \times \cdots \times [c_d,\infty)} e^{-\boldsymbol{\lambda} \cdot \boldsymbol{u}} I(\boldsymbol{\lambda}) \, \mathrm{d}\boldsymbol{\lambda}.$$
(20)

Theorem 2 (Properties of the Borel-Laplace transforms). Fix $\delta \in (0, \pi/2)$ and $\mathbf{c} \in \mathbb{R}^d_{>0}$.

- 1. For all $H \in \mathcal{P}(\mathbf{\Omega}_{\delta})$, $\mathcal{B}[H]$ defines an analytic function on \mathbf{K}_{δ} independent of the value of $\delta' \in (0, \delta)$. Moreover, $\mathcal{B}[H] \in \mathcal{P}(\mathbf{K}_{\delta^{\circ}})$ for all $\delta^{\circ} \in (0, \delta)$.
- 2. For all $I \in \mathcal{P}(\mathbf{K}_{\delta})$, $\mathcal{L}_{\mathbf{c}}[I]$ defines an analytic function on Ω_{δ} . Moreover, for all $\delta^{\circ} \in (0, \delta)$, there exists M > 0 such that $\mathcal{B} \circ \mathcal{L}_{\mathbf{c}}[I]$ is well-defined and analytic on $\{\lambda \in \mathbf{K}_{\delta^{\circ}} : \forall j, |\lambda_j| > M\}$.
- 3. (\mathcal{L}_{c} is a right inverse of \mathcal{B}). For all $I \in \mathcal{P}(\mathbf{K}_{\delta})$, we have $\mathcal{B} \circ \mathcal{L}_{c}[I] = I$. For all $H \in \mathcal{P}(\mathbf{\Omega}_{\delta})$, there exists a demi-entire function $E_{c} : \mathbf{\Omega}_{\delta} \to \mathbb{C}$ such that $\mathcal{L}_{c} \circ \mathcal{B}[H] = H + E_{c}$.

Corollary 3. The scaling function $I(\boldsymbol{\lambda})$ in Theorem 1(2) is $(-\Theta, \boldsymbol{\theta})$ -homogeneous and in $\mathcal{P}(\boldsymbol{K}_{\delta^{\circ}})$ for all $\delta^{\circ} \in (0, \delta)$. It is identically zero if and only if H is demi-entire on Ω_{δ} .

Outline of the rest of the paper. Sections 2 and 3 give the proofs of Theorems 1' and 2, respectively. In Section 4, I first provide some additional results on the classes of multivariate functions mentioned above Theorem 1, and then discuss the background of this paper and its relations to previous works.

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2 Proof of the multivariate transfer theorems

Contours and domains. We start by defining some contours and domains useful in the proof. Recall the definition (1) of the Δ -domain Δ_{δ} . For given values of $n_0 \gg 1$ and $\theta_j > 0$, we define the contours C_j and C_j^O by Figure 1(a): C_j is a contour inside Δ_{δ} and close to its boundary, while C_j^O is a portion of C_j close to the point z = 1. We denote by $\epsilon = \epsilon(\delta)$ the maximal distance from a point of C_j^O to 1. Let V_j^O be the image of C_j^O under the mapping $z \mapsto n_0^{\theta_j}(1-z)$, and V be the limit of V_j^O as $n_0 \to \infty$. We have $V_j^O \subset V \subset \Omega_{\delta}$, see Figure 1(b). The multivariate versions of these objects are defined by

$$\mathcal{C} = \mathcal{C}_1 \times \cdots \times \mathcal{C}_d \qquad \qquad \mathcal{C}^O = \mathcal{C}_1^O \times \cdots \times \mathcal{C}_d^O \\ \mathbf{V} = V \times \cdots \times V \qquad \qquad \mathbf{V}^O = V_1^O \times \cdots \times V_d^O .$$

Uniform growth bounds. Now let us derive some uniform bounds of functions on the above contours and domains. Recall that in the θ -diagonal limit, the vector $\boldsymbol{\lambda} = n_0^{-\theta} \boldsymbol{n} \equiv (n_0^{-\theta_1} n_1, \dots, n_0^{-\theta_d} n_d)$ remains in some compact subset of $(0, \infty)^d$ as $n_0 \to \infty$. Let $\lambda_{\min}, \lambda_{\max} > 0$ be such that $\lambda_j \in [\lambda_{\min}, \lambda_{\max}]$ for all $1 \le j \le d$. Let $\theta_{\min} = \min\{\theta_1, \dots, \theta_d\}$ and $\theta_{\max} = \max\{\theta_1, \dots, \theta_d\}$. We assume $n_0 \ge 2$.

We assume without loss of generality that δ , and therefore ϵ , is small enough so that all the local assumptions near $\boldsymbol{z} = \boldsymbol{1}$ in Theorem 1' hold globally in $\boldsymbol{B}_{1,\epsilon} := \{ \boldsymbol{z} \in \mathbb{C}^d \mid \forall 1 \leq j \leq d, |z_j - 1| < \epsilon \}$. In other words, the functions A_{reg} and $\boldsymbol{z} \mapsto H_k(\boldsymbol{1} - \boldsymbol{z})$ are well-defined on $\boldsymbol{\Delta}_{\delta} \cap \boldsymbol{B}_{1,\epsilon}$ and satisfy

$$\forall \boldsymbol{z} \in \boldsymbol{\Delta}_{\delta} \cap \boldsymbol{B}_{\boldsymbol{1},\epsilon}, \quad |F(\boldsymbol{z})| \leq C \cdot \left(|z_1 - 1|^{-M} + \dots + |z_d - 1|^{-M}\right), \tag{21}$$

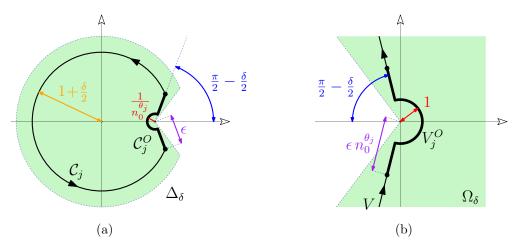


Figure 1: (a) The contour C_j is composed of two circle arcs and two line segments. The two arcs have their centers at 0 and 1, and their radii $1 + \delta/2$ and $1/n_0^{\theta_j}$, respectively. The line segments connect the two arcs and form an angle $\frac{\pi}{2} - \frac{\delta}{2}$ with the positive real axis at 1. The contour C_j^O is C_j with the large arc removed. (b) The contour V_j^O is the image of C_j^O under the mapping $z \mapsto n_0^{\theta_j}(1-z)$, and V is the limit of V_j^O as $n_0 \to \infty$ (which no longer depends on j). V is a special case of the curve $V_{\delta'}$ appeared in the definition (15) of $I(\lambda)$ in Theorem 1.

and there exist functions $A_{\operatorname{reg},j}$ analytic in $(\Delta_{\delta}^{j-1} \times \mathbb{C} \times \Delta_{\delta}^{d-j}) \cap B_{1,\epsilon}$ such that $A_{\operatorname{reg}} = A_{\operatorname{reg},1} + \cdots + A_{\operatorname{reg},d}$. Notice that $\mathcal{C}^O \subset \Delta_{\delta} \cap B_{1,\epsilon}$ and $|z_j - 1| \ge n_0^{-\theta_j}$ for all $z \in \mathcal{C}$. So the bound (21) implies that

$$\forall \boldsymbol{z} \in \boldsymbol{\mathcal{C}}^{O}, \quad |F(\boldsymbol{z})| \leq C \cdot \left(n_{0}^{M\theta_{1}} + \dots + n_{0}^{M\theta_{d}} \right) \leq dC \cdot n_{0}^{M\theta_{\max}}.$$
(22)

Since A is of polynomial type globally on Δ_{δ} , it satisfies (21) for all $z \in \Delta_{\delta}$ and (22) for all $z \in C$. We leave the reader to check the elementary fact that there exist $c, \mu > 0$ depending only on δ , θ_{\min} and λ_{\min} , λ_{\max} , such that

$$\forall z \in \mathcal{C}_j^O, \quad \left| z^{n_j + 1} \right| \ge c \cdot \exp\left(\mu \cdot n_0^{\theta_j} \Re(z - 1) \right).$$
(23)

And since $\mathfrak{Re}(z-1) \ge n_0^{-\theta_j}$ for $z \in \mathcal{C}_j^O$ and |z| > 1 for $z \in \mathcal{C}_j \setminus \mathcal{C}_j^O$, there exists $\tilde{c} = \tilde{c}(\delta, \theta_{\min}, \lambda_{\min}, \lambda_{\max}) > 0$ such that

$$\forall z \in \mathcal{C}_j, \quad \left| z^{n_j + 1} \right| \ge \tilde{c} \,. \tag{24}$$

Under the change of variable $u = n_0^{\theta_j}(1-z)$, the bound (23) becomes

$$\forall u \in V_j^O, \quad \left| (1 - n_0^{-\theta_j} u)^{n_j + 1} \right| \ge c \cdot e^{-\mu \cdot \mathfrak{Re}(u)} \,. \tag{25}$$

In Section 4, we will discuss several equivalent formulations of the polynomial-type condition, and their consequences on homogeneous functions. In particular, we will see in Lemma 7 that a homogeneous function of polynomial type *locally at* $\mathbf{0} \in \mathbf{\Omega}_{\delta}$ is always bounded by a Laurent polynomial globally on $\mathbf{\Omega}_{\delta}$. An easy consequence of Lemma 7 is the following global polynomial bound for the functions H_0, \ldots, H_m : there exists C, M > 0 such that

$$\forall \boldsymbol{u} \in \boldsymbol{V}, \quad |H_k(\boldsymbol{u})| \le C \cdot |u_1 \cdots u_d|^M \tag{26}$$

Analyticity of the homogeneous components. Before starting the proof of the main theorems, let us show the analyticity of the functions H_k ($0 \le k < m$) mentioned in the 7-th remark below Theorem 1'. Actually, we show that they are analytic on the larger domain Ω_{δ} , thanks to homogeneity.

Lemma 4. The functions H_0, \ldots, H_{m-1} in Theorem 1' have analytic continuations on Ω_{δ} .

Proof. Recall that by our assumption on the smallness of ϵ , the function A_{reg} is analytic on $\Delta_{\delta} \cap B_{1,\epsilon}$. Let $B_{0,\epsilon} = \{1 - z \mid z \in B_{1,\epsilon}\} \equiv \{z \in \mathbb{C}^d \mid \forall 1 \leq j \leq d, |z_j| < \epsilon\}$ and $A_{\text{sing}}(1 - z) = A(z) - A_{\text{reg}}(z)$ and $f_{\sigma}(\boldsymbol{u}) = \sigma^{-\theta_0^{(0)}} A_{\operatorname{sing}}(\sigma^{\theta} \boldsymbol{u})$. Then f_{σ} is analytic on $\Omega_{\delta} \cap \boldsymbol{B}_{0,\epsilon}$ for all $0 < \sigma \leq 1$, and the expansion (16) implies that $H_0(\boldsymbol{u}) = \lim_{\sigma \to 0^+} f_{\sigma}(\boldsymbol{u})$ for all $\boldsymbol{u} \in \Omega_{\delta} \cap \boldsymbol{B}_{0,\epsilon}$. Moreover, thanks to the polynomial-type bound, H_0, \ldots, H_m are bounded on all compact subsets of $\Omega_{\delta} \cap \boldsymbol{B}_{0,\epsilon}$. It follows that the family $(f_{\sigma})_{\sigma \in (0,1]}$ is uniformly bounded there. By Vitali's theorem, the convergence $f_{\sigma}(\boldsymbol{u}) \xrightarrow[\sigma \to 0]{} H_0(\boldsymbol{u})$ is uniform on compact subsets of $\Omega_{\delta} \cap \boldsymbol{B}_{0,\epsilon}$. Therefore H_0 is analytic on $\Omega_{\delta} \cap \boldsymbol{B}_{0,\epsilon}$.

The same argument applied to $A_{\text{sing}} - H_0$, $A_{\text{sing}} - H_0 - H_1$, etc. shows that H_1, \ldots, H_{m-1} are also analytic on $\Omega_{\delta} \cap B_{0,\epsilon}$. The homogeneity property $H_k(\sigma^{\theta} u) = \sigma^{\theta_0^{(k)}} H_k(u)$ extends this to all $u \in \Omega_{\delta}$. \Box

Proof of Theorem 1'. Now let us prove Theorem 1'. The proof will also imply Theorem 1 by considering cases in which only one term on the right hand side of (16) is nonzero. We follow the same general steps as the proof of the univariate transfer theorem in [5].

Step 1. Contour deformation and localization. The coefficient $[z^n]A(z)$ are related to A(z) by the Cauchy integral formula

$$[\boldsymbol{z}^{\boldsymbol{n}}]A(\boldsymbol{z}) = \left(\frac{1}{2\pi i}\right)^{d} \oint_{\boldsymbol{T}_{r}} \frac{A(\boldsymbol{z})}{\boldsymbol{z}^{\boldsymbol{n}+1}} \mathrm{d}\boldsymbol{z}$$
(27)

where $\mathbf{T}_r = \{ \mathbf{z} \in \mathbb{C}^d \mid \forall j, |z_j| = r \}$ is the polytorus of radius r for some r > 0 small enough. Thanks to the analyticity of $A(\mathbf{z})$ in $\mathbf{\Delta}_{\delta}$, we can deform the contour in (27) from \mathbf{T}_r to \mathbf{C} . By splitting the integral according to the partition $\mathbf{C} = \mathbf{C}^O \cup (\mathbf{C} \setminus \mathbf{C}^O)$, we obtain

$$[\boldsymbol{z}^{\boldsymbol{n}}]A(\boldsymbol{z}) = I_{\text{loc}} + \left(\frac{1}{2\pi i}\right)^d \int_{\boldsymbol{\mathcal{C}}\backslash\boldsymbol{\mathcal{C}}^O} \frac{A(\boldsymbol{z})}{\boldsymbol{z}^{\boldsymbol{n}+1}} d\boldsymbol{z} \quad \text{where} \quad I_{\text{loc}} = \left(\frac{1}{2\pi i}\right)^d \int_{\boldsymbol{\mathcal{C}}^O} \frac{A(\boldsymbol{z})}{\boldsymbol{z}^{\boldsymbol{n}+1}} d\boldsymbol{z} \,. \tag{28}$$

By the definition of \mathcal{C}^O , for each $z \in \mathcal{C} \setminus \mathcal{C}^O$, there is at least one $j^* \in \{1, \ldots, d\}$ such that $|z_{j^*}| = 1 + \delta/2$. Together with the growth bounds (22) and (24), this implies

$$\forall \boldsymbol{z} \in \boldsymbol{\mathcal{C}} \setminus \boldsymbol{\mathcal{C}}^{O}, \quad \left| \frac{A(\boldsymbol{z})}{\boldsymbol{z}^{\boldsymbol{n}+1}} \right| \leq \frac{dC \cdot n_{0}^{M\theta_{\max}}}{\tilde{c}^{d-1} \cdot (1+\delta/2)^{n_{j^{*}}+1}} \leq \frac{\tilde{C} \cdot n_{0}^{M\theta_{\max}}}{(1+\delta/2)^{\lambda_{\min} \cdot n_{0}^{\theta_{\min}}}} \leq C \cdot r^{n_{0}^{\theta_{\min}}}$$
(29)

for some constants r < 1 and C independent of n_0 . It follows that

$$[\boldsymbol{z}^{\boldsymbol{n}}]A(\boldsymbol{z}) = I_{\text{loc}} + O(r^{n_0^{\theta_{\min}}})$$
(30)

for some constant $r \in (0, 1)$ as $n_0 \to \infty$.

Step 2. Removing the demi-analytic term. Let $A_{sing}(1-z) = A(z) - A_{reg}(z)$. Then we have

$$I_{\text{loc}} = I_{\text{reg}} + I_{\text{sing}} := \left(\frac{1}{2\pi i}\right)^d \int_{\mathcal{C}^O} \frac{A_{\text{reg}}(\boldsymbol{z})}{\boldsymbol{z}^{\boldsymbol{n}+1}} d\boldsymbol{z} + \left(\frac{1}{2\pi i}\right)^d \int_{\mathcal{C}^O} \frac{A_{\text{sing}}(\boldsymbol{1}-\boldsymbol{z})}{\boldsymbol{z}^{\boldsymbol{n}+1}} d\boldsymbol{z}$$
(31)

Recall that the demi-analytic term A_{reg} admits a decomposition $A_{\text{reg}} = A_{\text{reg},1} + \cdots + A_{\text{reg},d}$, where each $A_{\text{reg},j}$ is analytic in $(\Delta_{\delta}^{j-1} \times \mathbb{C} \times \Delta_{\delta}^{d-j}) \cap B_{1,\epsilon}$. Within this domain, we can deform the *j*-th component the contour $\mathcal{C}^O \equiv \mathcal{C}_1^O \times \cdots \times \mathcal{C}_d^O$ to the arc $\tilde{\mathcal{C}}_j^O$ on the circle $\{z \in \mathbb{C} : |z| = 1 + \delta/2\}$ with the same endpoints as \mathcal{C}_j^O . Then, with the same reasoning as in Step 1, we obtain

$$\int_{\mathcal{C}^O} \frac{A_{\operatorname{reg},j}(\boldsymbol{z})}{\boldsymbol{z}^{\boldsymbol{n}+1}} \mathrm{d}\boldsymbol{z} = \int_{\mathcal{C}_1^O \times \dots \times \tilde{\mathcal{C}}_j^O \times \dots \times \mathcal{C}_d^O} \frac{A_{\operatorname{reg},j}(\boldsymbol{z})}{\boldsymbol{z}^{\boldsymbol{n}+1}} \mathrm{d}\boldsymbol{z} = O(r^{n_0^{\theta_{\min}}})$$
(32)

for some $r \in (0,1)$. Summing over j gives $I_{\text{reg}} = O(r^{n_0^{\theta_{\min}}})$. Together with Step 1, this implies

$$[\boldsymbol{z}^{\boldsymbol{n}}]A(\boldsymbol{z}) = I_{\text{sing}} + O(r^{n_0^{\theta_{\min}}})$$
(33)

for some constant $r \in (0,1)$ as $n_0 \to \infty$. In particular, when $A_{sing} = 0$, we obtain Case 1 of Theorem 1.

Step 3. Transfer. By linearity, the expansion (16) implies $I_{sing} = I_0 + \cdots + I_m$, where

$$I_{k} = \left(\frac{1}{2\pi i}\right)^{d} \int_{\mathcal{C}^{O}} \frac{H_{k}(\mathbf{1}-\mathbf{z})}{\mathbf{z}^{n+1}} d\mathbf{z} \quad \text{for } 0 \le k < m \qquad \text{and} \qquad I_{m} = \left(\frac{1}{2\pi i}\right)^{d} \int_{\mathcal{C}^{O}} \frac{o(H_{m}(\mathbf{1}-\mathbf{z}))}{\mathbf{z}^{n+1}} d\mathbf{z} \quad (34)$$

Thanks to the analyticity of H_0, \ldots, H_{m-1} in Lemma 4, we can apply Step 1 to them to obtain

$$[\boldsymbol{z}^{\boldsymbol{n}}]H_k(\boldsymbol{1}-\boldsymbol{z}) = I_k + O(r^{n_0^{\theta_{\min}}}) \quad \text{for all } 0 \le k < m.$$
(35)

Now let us show that $I_m = o(n_0^{-\Theta_m})$. To make use of the little-o estimate in the integrand, we want to further localize the contour \mathcal{C}^O : For a function $\varepsilon(x) \xrightarrow[x \to \infty]{} 0$, we define the little-o versions of the contours \mathcal{C}^O_i and \mathcal{C}^O by

$$\mathcal{C}_{j}^{o} = \{ z \in \mathcal{C}_{j}^{O} : |z - 1| < \varepsilon(n_{0}) \} \quad \text{and} \quad \mathcal{C}^{o} = \mathcal{C}_{1}^{o} \times \dots \times \mathcal{C}_{d}^{o}.$$
(36)

Contrary to \mathcal{C}^{O} , which has a small but fixed size ϵ , the contour \mathcal{C}^{O} shrinks to the point **1** when $n_{0} \to \infty$. According to the above definition, for all $\mathbf{z} \in \mathcal{C}^{O} \setminus \mathcal{C}^{O}$, there exists at least one $j^{*} \in \{1, \ldots, d\}$ such that $|z_{j^{*}}| \geq 1 + \sin(\delta/2) \cdot \varepsilon(n_{0})$. Then, similarly to (29), we have

$$\left| \forall \boldsymbol{\mathcal{C}}^{O} \setminus \boldsymbol{\mathcal{C}}^{o}, \quad \left| \frac{H_{m}(1-\boldsymbol{z})}{\boldsymbol{z}^{\boldsymbol{n}+1}} \right| \leq \frac{\tilde{C}n_{0}^{M\theta_{\max}}}{\left(1+\sin(\delta/2)\cdot\varepsilon(n_{0})\right)^{\lambda_{\min}\cdot n_{0}^{\theta_{\min}}}} \leq \mathsf{C} \cdot r^{\varepsilon(n_{0})\cdot n_{0}^{\theta_{\min}}}$$
(37)

for some r < 1 and C > 0. We choose a function $\varepsilon(x)$ such that $\varepsilon(x) \gg x^{-\theta_{\min}} \log(x)$ when $x \to \infty$. (For example, $\varepsilon(x) = x^{\theta_{\min}/2}$.) Then the right hand side of the above display decays faster than any negative power of n_0 when $n_0 \to \infty$. In particular, this implies

$$\int_{\mathcal{C}^{O} \setminus \mathcal{C}^{o}} \frac{o(H_{m}(\mathbf{1}-\boldsymbol{z}))}{\boldsymbol{z}^{n+1}} \mathrm{d}\boldsymbol{z} = o\left(n_{0}^{-\Theta_{m}}\right).$$
(38)

as $n_0 \to \infty$. On the other hand, since the contour \mathcal{C}^o shrinks to $\mathbf{z} = \mathbf{1}$ when $n_0 \to \infty$, for any $\tau > 0$, there exists n_0^* such that

$$\forall \boldsymbol{z} \in \boldsymbol{\mathcal{C}}^{o}, \quad |o(H_m(\boldsymbol{1}-\boldsymbol{z}))| \leq \tau |H_m(\boldsymbol{1}-\boldsymbol{z})|$$
(39)

for all $n_0 \ge n_0^*$. By plugging this inequality into the integral over \mathcal{C}^o and making the change of variable $u = n_0^{\theta}(1-z) \equiv (n_0^{\theta_1}(1-z_1), \ldots, n_0^{\theta_d}(1-z_d))$, we obtain

$$\left| \int_{\mathcal{C}^{o}} \frac{o(H_{m}(\mathbf{1}-\boldsymbol{z}))}{\boldsymbol{z}^{n+1}} \mathrm{d}\boldsymbol{z} \right| \leq \tau \int_{\mathbf{V}^{o}} \left| \frac{H_{m}(n_{0}^{-\boldsymbol{\theta}}\boldsymbol{u})}{\left(\mathbf{1}-n_{0}^{-\boldsymbol{\theta}}\boldsymbol{u}\right)^{n+1}} \right| \frac{\mathrm{d}\boldsymbol{u}}{n_{0}^{\theta_{1}+\dots+\theta_{d}}} = \frac{\tau}{n_{0}^{\Theta_{m}}} \int_{\mathbf{V}^{o}} \left| \frac{H_{m}(\boldsymbol{u})}{\left(\mathbf{1}-n_{0}^{-\boldsymbol{\theta}}\boldsymbol{u}\right)^{n+1}} \right| |\mathrm{d}\boldsymbol{u}| \quad (40)$$

for $n_0 \ge n_0^*$, where V^o is the image of \mathcal{C}^o under the change of variable, and the last equality used the $(\theta_0^{(m)}, \theta)$ -homogeneity of H_m . The bounds (25) and (26) imply that the integrand on the right hand side is bounded by

$$\frac{C|u_1\cdots u_d|^M}{c^d e^{-\mu\cdot\Re\mathfrak{e}(u_1+\cdots+u_d)}} = c^{-d}C\prod_{j=1}^d |u_j|^M e^{\mu\Re\mathfrak{e}(u_j)},$$
(41)

which is integrable on $V \supset V^o$. Hence the inegral on V^o is bounded by a constant independent of n_0 . It follows that $\int_{\mathcal{C}^o} \frac{o(H_m(1-z))}{z^{n+1}} dz = o(n^{-\Theta_m})$, Adding this to (38) gives $I_m = o(n_0^{-\Theta_m})$. Together with the conclusions of Steps 2 and 3, this implies the asymptotic expansion (17) in Theorem 1'.

When the little-o estimate in (16) is replaced by a big-O, the same proof (actually simpler, since one no longer needs to localize \mathcal{C}^O to \mathcal{C}^o) shows the expansion (17) with $o(n_0^{-\Theta_m})$ replaced by $O(n_0^{-\Theta_m})$.

Step 4. Coefficient asymptotics of homogeneous functions. It remains to prove the asymptotic expansion (14) for a general (θ_0, θ) -homogeneous function H satisfying the assumptions of Theorem 1. With the same argument as in Step 3, there exists $r \in (0, 1)$ such that

$$[\boldsymbol{z}^{\boldsymbol{n}}]H(\boldsymbol{1}-\boldsymbol{z}) = J + O(r^{n_0^{\theta_{\min}}}) \quad \text{with} \quad J := \left(\frac{1}{2\pi i}\right)^d \int_{\boldsymbol{\mathcal{C}}^O} \frac{H(\boldsymbol{1}-\boldsymbol{z})}{\boldsymbol{z}^{\boldsymbol{n}+1}} \mathrm{d}\boldsymbol{z}$$
(42)

We perform the same change of variable as in Step 3, which gives

$$J = \left(\frac{1}{2\pi i}\right)^d \int_{\mathbf{V}^O} \frac{H(n_0^{-\boldsymbol{\theta}}\boldsymbol{u})}{(1 - n_0^{-\boldsymbol{\theta}}\boldsymbol{u})^{\boldsymbol{n}+1}} \frac{\mathrm{d}\boldsymbol{u}}{n_0^{\theta_1 + \dots + \theta_d}} = \frac{1}{n_0^{\Theta}} \cdot \left(\frac{1}{2\pi i}\right)^d \int_{\mathbf{V}^O} \frac{H(\boldsymbol{u})}{(1 - n_0^{-\boldsymbol{\theta}}\boldsymbol{u})^{\boldsymbol{n}+1}} \mathrm{d}\boldsymbol{u}$$
(43)

where $\Theta = \theta_0 + \cdots + \theta_d$ as defined in Theorem 1. According to the definition (3) of the coefficients $g_{k,l}$, we have

$$\left(1 - n_0^{-\theta}u\right)^{-\lambda n_0^{\theta} - 1} = e^{\lambda u} G\left(n_0^{-\theta}u, \lambda u\right) = e^{\lambda u} \sum_{k=0}^{\infty} \left(\sum_{l=0}^k g_{k,l} \lambda^l u^{k+l}\right) \frac{1}{n_0^{k\cdot\theta}}.$$
(44)

Applying the above formula to each factor of the product $\prod_{j=1}^{d} \left(1 - n_0^{-\theta_j} u_j\right)^{-n_j-1}$ gives

$$\frac{1}{(1-n_0^{-\theta}\boldsymbol{u})^{n+1}} = e^{\boldsymbol{\lambda}\cdot\boldsymbol{u}} \sum_{\boldsymbol{k}\in\mathbb{N}^d} \left(\sum_{\boldsymbol{l}\leq\boldsymbol{k}} g_{\boldsymbol{k},\boldsymbol{l}} \,\boldsymbol{\lambda}^{\boldsymbol{l}} \boldsymbol{u}^{\boldsymbol{k}+\boldsymbol{l}}\right) \frac{1}{n_0^{\boldsymbol{k}\cdot\boldsymbol{\theta}}} \,. \tag{45}$$

If we treat the right hand side as a formal sum over k, and approximate the contour V^O by its limit V, then (43) would imply heuristically that:

$$J \approx \frac{1}{n^{\Theta}} \sum_{\boldsymbol{k} \in \mathbb{N}^{d}} \frac{1}{n_{0}^{\boldsymbol{k} \cdot \boldsymbol{\theta}}} \cdot \left(\frac{1}{2\pi i}\right)^{d} \int_{\boldsymbol{V}} \left(\sum_{\boldsymbol{l} \leq \boldsymbol{k}} g_{\boldsymbol{k}, \boldsymbol{l}} \boldsymbol{\lambda}^{\boldsymbol{l}} \boldsymbol{u}^{\boldsymbol{k} + \boldsymbol{l}}\right) e^{\boldsymbol{\lambda} \cdot \boldsymbol{u}} H(\boldsymbol{u}) d\boldsymbol{u}$$
$$= \frac{1}{n^{\Theta}} \sum_{\boldsymbol{k} \in \mathbb{N}^{d}} \frac{1}{n_{0}^{\boldsymbol{k} \cdot \boldsymbol{\theta}}} \cdot \sum_{\boldsymbol{l} \leq \boldsymbol{k}} g_{\boldsymbol{k}, \boldsymbol{l}} \boldsymbol{\lambda}^{\boldsymbol{l}} \boldsymbol{\partial}_{\boldsymbol{\lambda}}^{\boldsymbol{k} + \boldsymbol{l}} \left(\left(\frac{1}{2\pi i}\right)^{d} \int_{\boldsymbol{V}} e^{\boldsymbol{\lambda} \cdot \boldsymbol{u}} H(\boldsymbol{u}) d\boldsymbol{u}\right)$$
$$= \frac{1}{n^{\Theta}} \sum_{\boldsymbol{k} \in \mathbb{N}^{d}} \frac{1}{n_{0}^{\boldsymbol{k} \cdot \boldsymbol{\theta}}} \cdot D_{\boldsymbol{k}} I(\boldsymbol{\lambda}).$$

where the function $I(\lambda)$ and the differential operator D_k are defined as in Theorem 1.

Now let us show that the last line is indeed an asymptotic expansion of J. In other words, for any N > 0, we have

$$n_0^{\Theta} \cdot J = \sum_{\boldsymbol{k} \cdot \boldsymbol{\theta} < N} \frac{D_{\boldsymbol{k}} I(\boldsymbol{\lambda})}{n_0^{\boldsymbol{k} \cdot \boldsymbol{\theta}}} + O\left(\frac{1}{n_0^N}\right)$$
(46)

when $n_0 \to \infty$. For this, let us go back to (44). It can be seen as the Taylor series expansion of the function $f_y(x) = (1-x)^{-y/x-1}$ around x = 0 evaluated at $x = n_0^{-\theta} u$ and $y = \lambda u$. The corresponding Taylor expansion with remainder term writes

$$f_y(x) = e^y \sum_{k=0}^m \left(\sum_{l=0}^k g_{k,l} y^l\right) x^k + R_m(x,y) \quad \text{with} \quad R_m(x,y) = \frac{1}{m!} \int_0^x (x-\xi)^m f_y^{(m+1)}(\xi) \mathrm{d}\xi.$$
(47)

It is not hard to show by induction that there are functions $\varphi_{k,m}$ continuous on the unit disk, such that

$$f_{y}^{(m)}(x) = f_{y}(x) \cdot \sum_{k=0}^{m} \varphi_{k,m}(x) y^{k}, \qquad (48)$$

It follows that for each m, there exists a constant $C_m > 0$ such that $|f_y^{(m)}(x)| \leq C_m(1+|y|^m) \cdot |f_y(x)|$ for all $|x| \leq 1/2$ and $y \in \mathbb{C}$. Plugging this into the definition of $R_m(x,y)$ gives that

$$\forall |x| \le 1/2, \, \forall y \in \mathbb{C}, \quad |R_m(x,y)| \le \frac{C_m}{(m+1)!} (1+|y|^m) \cdot |x|^{m+1} \sup_{\xi \in [0,x]} \left| f_y(\xi) \right| \tag{49}$$

Consider the case where $(x, y) = (n_0^{-\theta_j} u, \lambda u)$ with $n_0 > 0, \lambda \in [\lambda_{\min}, \lambda_{\max}]$ and $u \in V_j^O$. In this case the condition $|x| \leq 1/2$ is satisfied because $|u| \leq \epsilon n_0^{\theta_j}$ on V_j^O and $\epsilon = \epsilon(\delta)$ is assumed to be small enough. Since $|u| \geq 1$ on V_j^O , the term $(1 + |y|^m)$ is bounded by a constant times $|u|^m$ uniformly for $u \in V_j^O$. Moreover, for all $\xi \in [0, x]$, there exists $\tilde{n}_0 > 0$ such that $\xi = \tilde{n}_0^{-\theta_j} u$. Thus we can apply the bound (25) to see that $|f_y(\xi)| = |(1 - \tilde{n}_0^{-\theta_j} u)^{-\lambda u - 1}| \leq c^{-1} e^{\mu \cdot \Re \epsilon(u)}$ for all $\xi \in [0, x]$. It follows that

$$\forall u \in V_j^O, \quad \left| R_m(n_0^{-\theta_j}u, \lambda u) \right| \le \tilde{C}_m \frac{|u|^{2m+1} e^{\mu \cdot \mathfrak{Re}(u)}}{n_0^{(m+1)\theta_j}} \tag{50}$$

for some constant \tilde{C}_m depending only on δ , λ_{\min} , λ_{\max} and m. Then, the Taylor expansion (47) becomes

$$(1 - n_0^{-\theta_j} u)^{-\lambda n_0^{\theta_j} - 1} = e^{\lambda u} \sum_{k=0}^m \left(\sum_{l=0}^k g_{k,l} \,\lambda^l u^{k+l} \right) \frac{1}{n_0^{k \cdot \theta_j}} + O\left(\frac{|u|^{2m+1} e^{\mu \cdot \mathfrak{Re}(u)}}{n_0^{(m+1)\theta_j}} \right) \,, \tag{51}$$

where the big-O estimate is uniform with respect to $u \in V_j^O$ as $n_0 \to \infty$. Now take $m > N/\theta_{\min}$ and replace each factor in $\prod_{j=1}^d (1 - n_0^{-\theta_j} u_j)^{-n_j-1}$ by the right hand side of the above formula. After expanding the resulting product, we obtain a *finite* sum. By collecting all the terms of order $O(1/n_0^N)$ together, we obtain an expansion of the form

$$\frac{1}{(\mathbf{1} - n_0^{-\boldsymbol{\theta}} \boldsymbol{u})^{n+1}} = e^{\boldsymbol{\lambda} \cdot \boldsymbol{u}} \sum_{\boldsymbol{k} \cdot \boldsymbol{\theta} < N} \left(\sum_{\boldsymbol{l} \le \boldsymbol{k}} g_{\boldsymbol{k}, \boldsymbol{l}} \, \boldsymbol{\lambda}^{\boldsymbol{l}} \boldsymbol{u}^{\boldsymbol{k}+\boldsymbol{l}} \right) \frac{1}{n_0^{\boldsymbol{k} \cdot \boldsymbol{\theta}}} + O\left(\frac{R(\boldsymbol{u})}{n_0^N}\right), \tag{52}$$

where the big-O estimate is uniform on $\boldsymbol{u} \in \boldsymbol{V}^O$. By following the above calculation more closely, it is not hard to see that we can choose $R(\boldsymbol{u}) = |u_1 \cdots u_d|^{\tilde{M}} e^{\boldsymbol{\mu} \cdot \mathfrak{Re}(u_1 + \cdots + u_d)}$ for some $\boldsymbol{\mu}, \tilde{M} > 0$. With the bound (26) for $H(\boldsymbol{u})$, it follows that $H(\boldsymbol{u})R(\boldsymbol{u})$ is integrable on $\boldsymbol{V} \supset \boldsymbol{V}^O$. So we can integrable the previous display term by term with respect to $\boldsymbol{u} \in \boldsymbol{V}^O$ to obtain

$$\int_{\mathbf{V}^O} \frac{H(\boldsymbol{u})}{(\mathbf{1} - n_0^{-\boldsymbol{\theta}} \boldsymbol{u})^{\boldsymbol{n}+1}} d\boldsymbol{u} = \sum_{\boldsymbol{k} \cdot \boldsymbol{\theta} < N} \int_{\mathbf{V}^O} H(\boldsymbol{u}) e^{\boldsymbol{\lambda} \cdot \boldsymbol{u}} \left(\sum_{\boldsymbol{l} \le \boldsymbol{k}} g_{\boldsymbol{k}, \boldsymbol{l}} \, \boldsymbol{\lambda}^{\boldsymbol{l}} \boldsymbol{u}^{\boldsymbol{k}+\boldsymbol{l}} \right) d\boldsymbol{u} \cdot \frac{1}{n_0^{\boldsymbol{k} \cdot \boldsymbol{\theta}}} + O\left(\frac{1}{n_0^N}\right).$$
(53)

Like in Step 1, replacing the contour \mathbf{V}^O by \mathbf{V} in each term of the above equation only produces an error of order $O(r^{n_0^{\theta_{\min}}})$ as $n_0 \to \infty$. The resulting equation divided by $(2\pi i)^d$ gives exactly (46).

The expansion (14) follows readily from (42) and (46). This concludes the proof of Theorem 1'. \Box

3 The Borel-Laplace transforms

In this section, we prove the three statements in Theorem 2.

Proof of Theorem 2(1). Fix $\delta \in (0, \frac{\pi}{2})$ and $H \in \mathcal{P}(\Omega_{\delta})$. Recall that $\mathcal{B}[H]$ is defined as the integral

$$I_{\boldsymbol{V}_{\delta'}}(\boldsymbol{\lambda}) := \left(\frac{1}{2\pi i}\right)^d \int_{\boldsymbol{V}_{\delta'}} e^{\boldsymbol{\lambda} \cdot \boldsymbol{u}} H(\boldsymbol{u}) \mathrm{d}\boldsymbol{u},\tag{54}$$

on a contour $V_{\delta'}$ chosen from the class of contours

 $\boldsymbol{\mathcal{V}}_{\delta'} := \left\{ V_1 \times \dots \times V_d \mid \forall 1 \le j \le d, \, V_j \text{ is a piecewise smooth curve in } \Omega_{\delta} \\ \text{which coincide with } \partial \Omega_{\delta'} \text{ outside a bounded set, as in Figure 2(a)} \right\}.$

Thanks to the analyticity of H on the domain $\Omega_{\delta} = (\Omega_{\delta})^d$, it is clear that for each fixed δ' , the value of the integral $I_{\mathbf{V}_{\delta'}}(\boldsymbol{\lambda})$ does not depend on the choice of the contour $\mathbf{V}_{\delta'}$ within the class $\boldsymbol{\mathcal{V}}_{\delta'}$.

First, let us fix $\delta' \in (0, \delta)$ and $V_{\delta'} \in \mathcal{V}_{\delta'}$ and show that the integral $I_{V_{\delta'}}(\lambda)$ is absolutely convergent and analytic with respect to $\lambda \in K_{\delta'}$. Let $U_{\delta',r}$ be the closure of $(\mathbb{C} \setminus \Omega_{\delta'}) \cup B_{0,r}$, where $B_{0,r} \subset \mathbb{C}$ is the disk of radius r around the origin. It is a simple exercise to show that for $\delta^{\circ} \in (0, \delta')$ and r > 0, we have

$$\forall u \in U_{\delta',r}, \ \forall \lambda \in K_{\delta^{\circ}}, \quad \left| e^{\lambda u} \right| \leq C_*^{r|\lambda|} e^{-\sigma_*|\lambda| \cdot |u|}, \tag{55}$$

with $\sigma_* = \sin(\delta' - \delta^\circ) > 0$ and $C_* = e^{1+\sigma_*}$. Since each component of the contour $V_{\delta'} \in \mathcal{V}_{\delta'}$ is bounded away from the origin and H is of polynomial type, there exist $\tilde{C}, M > 0$ such that

$$\forall \boldsymbol{u} \in \boldsymbol{V}_{\delta'}, \quad |H(\boldsymbol{u})| \leq \tilde{C} \cdot \left(|u_1|^M + \dots + |u_d|^M\right) \,. \tag{56}$$

On the other hand, there exists r > 0 such that $V_{\delta'} \subset U_{\delta',r}$. So the bound (55) implies that for any bounded set $S \subseteq K_{\delta^{\circ}}$, there exists $C, \sigma > 0$, such that

$$\forall \boldsymbol{u} \in \boldsymbol{V}_{\delta'}, \ \forall \boldsymbol{\lambda} \in \boldsymbol{S}, \quad \left| e^{\boldsymbol{\lambda} \cdot \boldsymbol{u}} \right| \le C \cdot e^{-\sigma(|u_1| + \dots + |u_d|)}.$$
(57)

Up to increasing the value of C and decreasing σ , the polynomial bound on the right hand side of (56) can be absorbed by the exponential decay in the above display. Therefore we have

$$\forall \boldsymbol{u} \in \boldsymbol{V}_{\delta'}, \ \forall \boldsymbol{\lambda} \in \boldsymbol{S}, \quad \left| e^{\boldsymbol{\lambda} \cdot \boldsymbol{u}} H(\boldsymbol{u}) \right| \le C \cdot e^{-\sigma(|\boldsymbol{u}_1| + \dots + |\boldsymbol{u}_d|)}$$
(58)

for some $C, \sigma > 0$. Since the right hand side is independent of $\lambda \in S$ and integrable on $V_{\delta'}$, it follows that the integral $I_{V_{\delta'}}(\lambda)$ is absolutely convergent and analytic with respect to $\lambda \in S$. And since this is true for all $\delta^{\circ} \in (0, \delta')$ and bounded set $S \subset K_{\delta^{\circ}}$, the integral $I_{V_{\delta'}}$ defines an analytic function on $K_{\delta'}$.

Now let us show that $I_{V_{\delta'}}(\lambda)$ is independent of $\delta' \in (0, \delta)$ as well. For this, we fix $\delta', \delta'' \in (0, \delta)$ and $\lambda \in \mathbf{K}_{\min(\delta', \delta'')}$. For each N > 0, let $V_{\delta''}^N$ be the contour obtained by deforming each component of $V_{\delta''}$ inside the disk $B_{0,N}$ to coincide with the corresponding component of $V_{\delta'}$ there, while keeping it unchanged outside $B_{0,N}$. See Figure 2(a). By the triangular inequality, we have

$$\left| \left(\int_{\boldsymbol{V}_{\delta'}} - \int_{\boldsymbol{V}_{\delta''}} \right) e^{\boldsymbol{\lambda} \cdot \boldsymbol{u}} H(\boldsymbol{u}) \mathrm{d}\boldsymbol{u} \right| = \left| \left(\int_{\boldsymbol{V}_{\delta'}} - \int_{\boldsymbol{V}_{\delta''}} \right) e^{\boldsymbol{\lambda} \cdot \boldsymbol{u}} H(\boldsymbol{u}) \mathrm{d}\boldsymbol{u} \right| \le \int_{\boldsymbol{V}_{\delta'}} \Delta \boldsymbol{V}_{\delta''}^{N} \left| e^{\boldsymbol{\lambda} \cdot \boldsymbol{u}} H(\boldsymbol{u}) \right| \cdot \left| \mathrm{d}\boldsymbol{u} \right| , \quad (59)$$

where $V_{\delta'} \Delta V_{\delta''}^N$ is the symmetric difference between $V_{\delta'}$ and $V_{\delta''}^N$. Let

$$\boldsymbol{S}_{n} = \left\{ \boldsymbol{u} \in \boldsymbol{V}_{\delta'} \,\Delta \, \boldsymbol{V}_{\delta''}^{N} \,\middle|\, n \leq \max_{j} |u_{j}| < n+1 \right\}.$$
(60)

It is not hard to see that $S_n = \emptyset$ for n < N, and the *d*-dimensional Lebesgue measure of S_n is bounded by cn^d for some constant $c = c(\delta', \delta'', d)$. Moreover, since each point $\lambda \in K_{\min(\delta', \delta'')}$ belongs to K_{δ° for some $\delta^\circ < \min(\delta', \delta'')$, we can deduce from (58) that $|e^{\lambda \cdot u}H(u)| \le Ce^{-\sigma n}$ for all $u \in S_n$ and all n. (Here we also use the observation that (58) remains valid when $V_{\delta'}$ is replaced by $V_{\delta'}^N$). It follows that

$$\int_{\boldsymbol{V}_{\delta'}\Delta\boldsymbol{V}_{\delta''}} \left| e^{\boldsymbol{\lambda}\cdot\boldsymbol{u}} H(\boldsymbol{u}) \right| \cdot |\mathrm{d}\boldsymbol{u}| \leq \sum_{n=N}^{\infty} C e^{-\sigma n} \cdot cn^{d} \xrightarrow[N \to \infty]{} 0.$$
(61)

Together with (59), this implies that $I_{V_{\delta'}}(\lambda) = I_{V_{\delta''}}(\lambda)$ for all $\lambda \in K_{\min(\delta',\delta'')}$. In this sense, $\mathcal{B}[H] = I_{V_{\delta'}}$ is independent of $\delta' \in (0, \delta)$ and thus defines an analytic function on K_{δ} .

To show that $\mathcal{B}[H] \in \mathcal{P}(\mathbf{K}_{\delta^{\circ}})$ for all $\delta^{\circ} \in (0, \delta)$, let us fix $\lambda \in \mathbf{K}_{\delta^{\circ}}$ and $\delta' \in (\delta^{\circ}, \delta)$, and consider the particular contour $\mathbf{V}_{\delta'} = V_1 \times \cdots \times V_d$ with $V_j = \partial U_{\delta', 1/|\lambda_j|}$, where $U_{\delta', r}$ is defined above (55). By the bound (55), we have

$$\forall \boldsymbol{u} \in \boldsymbol{V}_{\delta'}, \quad \left| e^{\boldsymbol{\lambda} \cdot \boldsymbol{u}} \right| \le C_*^d \cdot e^{-\sigma_*(|\lambda_1 u_1| + \dots + |\lambda_d u_d|)} \tag{62}$$

Since the contour V_j is not bounded away from the origin when $|\lambda_j| \to \infty$, we need to use the complete bound (11) for the function of polynomial type $H \in \mathcal{P}(\mathbf{\Omega}_{\delta})$ here, which implies

$$\forall \boldsymbol{u} \in \boldsymbol{V}_{\delta'}, \quad \left| e^{\boldsymbol{\lambda} \cdot \boldsymbol{u}} H(\boldsymbol{u}) \right| \le C \cdot \left(|u_1|^M + |u_1|^{-M} + \dots + |u_d|^M + |u_d|^{-M} \right) \cdot e^{-\sigma_* \cdot (|\lambda_1 u_1| + \dots + |\lambda_d u_d|)} \quad (63)$$

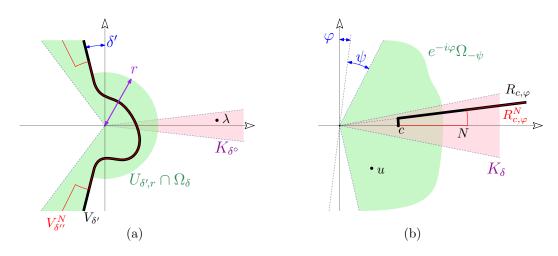


Figure 2: (a) Some domains and contours used in the proof of Theorem 2(1). (b) Some domains and contours used in the proof of Theorem 2(2).

for some C, M > 0. We integrate the above bound on $V_{\delta'}$ and make the change of variables $v_j = \lambda_j u_j$. Notice that the contour $\lambda V_{\delta'} \equiv (\lambda_1 V_1) \times \cdots \times (\lambda_d V_d)$ for the variable v after this change no longer depends on λ . This gives us

$$\begin{split} \left| \mathcal{B}[H](\boldsymbol{\lambda}) \right| &= \left| \left(\frac{1}{2\pi i} \right)^d \int_{\boldsymbol{V}_{\delta'}} e^{\boldsymbol{\lambda} \cdot \boldsymbol{u}} H(\boldsymbol{u}) \mathrm{d}\boldsymbol{u} \right| \\ &\leq \frac{C}{(2\pi)^d} \int_{\boldsymbol{\lambda} \boldsymbol{V}_{\delta'}} \left(\left| \frac{v_1}{\lambda_1} \right|^M + \left| \frac{v_1}{\lambda_1} \right|^{-M} + \dots + \left| \frac{v_d}{\lambda_d} \right|^M + \left| \frac{v_d}{\lambda_d} \right|^{-M} \right) \frac{e^{-\sigma_*(|v_1| + \dots + |v_d|)} |\mathrm{d}\boldsymbol{v}|}{|\lambda_1 \cdots \lambda_d|} \\ &\leq \frac{\max\left(|\lambda_1|^M, |\lambda_1|^{-M}, \dots, |\lambda_d|^M, |\lambda_d|^{-M} \right)}{|\lambda_1 \cdots \lambda_d|} \cdot J \end{split}$$

where the integral

$$J = \frac{C}{(2\pi)^d} \int_{\lambda V_{\delta'}} \left(|v_1|^M + |v_1|^{-M} + \dots + |v_d|^M + |v_d|^{-M} \right) \cdot e^{-\sigma_*(|v_1| + \dots + |v_d|)} |\mathrm{d}\boldsymbol{v}|$$
(64)

is finite and independent of λ . It follows that

$$\forall \boldsymbol{\lambda} \in \boldsymbol{K}_{\delta^{\circ}}, \quad \left| \mathcal{B}[H](\boldsymbol{\lambda}) \right| \le J \cdot \frac{|\lambda_1|^M + |\lambda_1|^{-M} + \dots + |\lambda_d|^M + |\lambda_d|^{-M}}{|\lambda_1 \cdots \lambda_d|}.$$
(65)

We will see in Lemma 8 that the monomial function $\lambda \mapsto \lambda^{\alpha}$ is in $\mathcal{P}(\mathbf{K}_{\delta^{\circ}})$, for all $\alpha \in \mathbb{R}^d$ and $\delta^{\circ} \in (0, \pi)$. This allows us to bound the right hand side of the above display by a polynomial-type bound of the form (11). Hence we have $\mathcal{B}[H] \in \mathcal{P}(\mathbf{K}_{\delta^{\circ}})$, for all $\delta^{\circ} \in (0, \delta)$.

Proof of Theorem 2(2). Fix $\boldsymbol{c} \in \mathbb{R}_{>0}^d$, $\delta \in (0, \frac{\pi}{2})$ and $I \in \mathcal{P}(\boldsymbol{K}_{\delta})$. For each $\boldsymbol{\varphi} \in (-\delta, \delta)^d$, we define $\boldsymbol{R}_{\boldsymbol{\varphi}} = R_{c_1,\varphi_1} \times \cdots \times R_{c_d,\varphi_d} \subset \boldsymbol{K}_{\delta}$, where $R_{c,\varphi} \subset \mathbb{C}$ is the contour which starts at $c \in \mathbb{R}_{>0}$, then follows an arc on the circle $\partial B_{0,|c|}$, and then goes to ∞ along the ray $\{\lambda \in \mathbb{C}_* | \arg(\lambda) = \varphi\}$, as in Figure 2(b). Let

$$H_{\varphi}(\boldsymbol{u}) = \int_{\boldsymbol{R}_{\varphi}} e^{-\boldsymbol{\lambda} \cdot \boldsymbol{u}} I(\boldsymbol{\lambda}) \mathrm{d}\boldsymbol{\lambda} \,. \tag{66}$$

Recall that $\Omega_{\psi} = (\Omega_{\psi})^d$ with $\Omega_{\psi} \equiv K_{\frac{\pi}{2}+\psi} := \left\{ u \in \mathbb{C}_* : |\arg(u)| < \frac{\pi}{2} + \psi \right\}$, for any $\psi \in (-\frac{\pi}{2}, \frac{\pi}{2})$.

First, let us show that the integral $H_{\varphi}(\boldsymbol{u})$ is absolutely convergent and analytic with respect to $\boldsymbol{u} \in e^{-i\varphi} \boldsymbol{\Omega}_0 := (e^{-i\varphi_1} \boldsymbol{\Omega}_0) \times \cdots \times (e^{-i\varphi_d} \boldsymbol{\Omega}_0)$. The proof is similar to the one for the integral $I_{\boldsymbol{V}_{\delta'}}(\boldsymbol{\lambda})$: It is a simple exercise to show that for any $\varphi \in (-\delta, \delta), \ \psi \in (0, \frac{\pi}{2})$ and $c \in \{c_1, \ldots, c_d\}$ we have

$$\forall \lambda \in R_{c,\varphi}, \ \forall u \in e^{-i\varphi}\Omega_{-\psi}, \quad \left|e^{-\lambda u}\right| \le C_*^{|u|}e^{-\sigma_*|\lambda| \cdot |u|} \tag{67}$$

with $\sigma_* = \sin(\psi) > 0$ and $C_* = e^{(1+\sigma_*)\max_j |c_j|}$. Following the same steps as for the bounds (56)–(58), we deduce that for any bounded subset $\mathbf{S} \subset e^{-i\varphi} \mathbf{\Omega}_{-\psi}$, where $\psi \in (0, \frac{\pi}{2})$ and $\varphi \in (-\delta, \delta)^d$, there exists $C, \sigma > 0$ such that

$$\forall \boldsymbol{\lambda} \in \boldsymbol{R}_{\boldsymbol{\varphi}}, \ \forall \boldsymbol{u} \in \boldsymbol{S}, \quad \left| e^{\boldsymbol{\lambda} \cdot \boldsymbol{u}} I(\boldsymbol{\lambda}) \right| \le C \cdot e^{-\sigma(|\lambda_1| + \dots + |\lambda_d|)}.$$
(68)

Since the right hand side is independent of $\boldsymbol{u} \in \boldsymbol{S}$ and integrable on \boldsymbol{R}_{φ} , it follows that the integral $H_{\varphi}(\boldsymbol{u})$ is absolutely convergent and analytic with respect to $\boldsymbol{u} \in \boldsymbol{S}$. And since this is true for all bounded $\boldsymbol{S} \subset e^{-i\varphi} \boldsymbol{\Omega}_{-\psi}$ with $\psi \in (0, \frac{\pi}{2})$, we conclude that H_{φ} defines an analytic function on $e^{-i\varphi} \boldsymbol{\Omega}_0$.

Now let us show that $H_{\varphi}(\boldsymbol{u}) = H_{\boldsymbol{0}}(\boldsymbol{u})$ for all $\varphi \in (-\delta, \delta)^d$ and $\boldsymbol{u} \in \Omega_{-\delta}$. Fix some $\varphi \in (-\delta, \delta)^d$. For any N > 0, let $\boldsymbol{R}_{\varphi}^N = \boldsymbol{R}_{c_1,\varphi_1}^N \times \cdots \times \boldsymbol{R}_{c_d,\varphi_d}^N$, where $\boldsymbol{R}_{c,\varphi}^N$ is the curve which coincides with the interval [c, N] inside the disk $B_{0,N}$, while staying on the ray $\{\lambda \in \mathbb{C}_* | \arg(\lambda) = \varphi\}$ outside $B_{0,N}$. See Figure 2(b). By construction, the contour $\boldsymbol{R}_{\varphi}^N$ coincides with $\boldsymbol{R}_{\boldsymbol{0}}$ in the polydisk $(B_{0,N})^d$. Like in the previous proof for $I_{\boldsymbol{V}_{\delta'}}(\boldsymbol{\lambda})$, the sets

$$\tilde{\boldsymbol{S}}_{n} = \left\{ \boldsymbol{\lambda} \in \boldsymbol{R}_{\boldsymbol{\varphi}}^{N} \, \Delta \, \boldsymbol{R}_{\boldsymbol{0}} \, \big| \, n \leq \max_{j} |\lambda_{j}| < n+1 \right\}$$
(69)

satisfy $\tilde{\mathbf{S}}_n = \emptyset$ for all n < N and $\int_{\tilde{\mathbf{S}}_n} |d\mathbf{\lambda}| \leq cn^d$ for some $c < \infty$ independent of n. Moreover, since each point $\mathbf{u} \in \mathbf{\Omega}_{-\delta}$ belongs to $\mathbf{\Omega}_{-\psi} \cap e^{-i\varphi}\mathbf{\Omega}_{-\psi}$ for some $\psi \in (0, \frac{\pi}{2})$, we can deduce from (68) that $|e^{-\mathbf{\lambda} \cdot \mathbf{u}}I(\mathbf{\lambda})| \leq Ce^{-\sigma n}$ for all $\mathbf{\lambda} \in \tilde{\mathbf{S}}_n$ and all n. (Here we also use the fact that (68) remains valid when \mathbf{R}_{φ} is replaced by \mathbf{R}_{φ}^N). It follows that

$$\begin{aligned} \left| H_{\varphi}(\boldsymbol{u}) - H_{\boldsymbol{0}}(\boldsymbol{u}) \right| &= \left| \left(\int_{\boldsymbol{R}_{\varphi}^{N}} - \int_{\boldsymbol{R}_{\boldsymbol{0}}} \right) e^{-\boldsymbol{\lambda} \cdot \boldsymbol{u}} I(\boldsymbol{\lambda}) \mathrm{d}\boldsymbol{\lambda} \right| \leq \int_{\boldsymbol{R}_{\varphi}^{N} \Delta \boldsymbol{R}_{\boldsymbol{0}}} \left| e^{-\boldsymbol{\lambda} \cdot \boldsymbol{u}} I(\boldsymbol{\lambda}) \right| |\mathrm{d}\boldsymbol{\lambda}| \\ &\leq \sum_{n=N}^{\infty} C e^{-\sigma n} \cdot cn^{d} \xrightarrow[N \to \infty]{} 0 \end{aligned}$$

Therefore, $H_{\varphi}(\boldsymbol{u}) = H_{\boldsymbol{0}}(\boldsymbol{u})$ for all $\varphi \in (-\delta, \delta)^d$ and $\boldsymbol{u} \in \boldsymbol{\Omega}_{-\delta}$. And since each H_{φ} is analytic on $e^{-i\varphi}\boldsymbol{\Omega}_0$, the function $\mathcal{L}_{\boldsymbol{c}}[I] \equiv H_{\boldsymbol{0}}$ has an analytic continuation on the union $\boldsymbol{\Omega}_{\delta} = \bigcup_{\boldsymbol{\varphi} \in (-\delta, \delta)^d} e^{-i\varphi}\boldsymbol{\Omega}_0$.

Finally, let us fix $\delta^{\circ} \in (0, \delta)$ and $\lambda \in \mathbf{K}_{\delta^{\circ}}$, and show that $\mathcal{B} \circ \mathcal{L}_{\mathbf{c}}[I]$ is well-defined and analytic at λ when the $|\lambda_j|$'s are large enough. For this, let us first prove that for all $\delta'' \in (0, \delta)$, there exists m > 0 such that

$$\forall \boldsymbol{u} \in \boldsymbol{\Omega}_{\delta''} \cap \tilde{\boldsymbol{B}}, \quad |\mathcal{L}_{\boldsymbol{c}}[I](\boldsymbol{u})| \le e^{m \cdot (|u_1| + \dots + |u_d|)}.$$

$$\tag{70}$$

where $\tilde{\boldsymbol{B}} := \{ \boldsymbol{u} \in \mathbb{C}^d \, | \, \forall j, \, |u_j| \ge 1 \}$. Indeed, for all $\boldsymbol{\varphi} \in (-\delta, \delta)^d$ and $\psi \in (0, \frac{\pi}{2}), \, (67)$ implies that

$$\forall \boldsymbol{u} \in e^{-i\boldsymbol{\varphi}} \boldsymbol{\Omega}_{-\psi}, \; \forall \boldsymbol{\lambda} \in \boldsymbol{R}_{\boldsymbol{\varphi}}, \quad \left| e^{-\boldsymbol{\lambda} \cdot \boldsymbol{u}} \right| \le C_*^{|u_1| + \dots + |u_d|} \cdot e^{-\sigma_*(|\lambda_1 u_1| + \dots + |\lambda_d u_d|)}. \tag{71}$$

In addition, since $I \in \mathcal{P}(\mathbf{K}_{\delta})$ and $\mathbf{R}_{\varphi} \subset \mathbf{K}_{\delta}$ is bounded away from **0**, there exist $\tilde{C}, M > 0$ such that

$$\forall \boldsymbol{\lambda} \in \boldsymbol{R}_{\boldsymbol{\varphi}}, \quad \left| I(\boldsymbol{\lambda}) \right| \leq \tilde{C} \left(|\lambda_1|^M + \dots + |\lambda_d|^M \right) \,. \tag{72}$$

For $\boldsymbol{u} \in \boldsymbol{B}$, we can increase the value of C_* and decrease that of σ_* on the right hand side of (71) to absorbe the polynomial function on right hand side of (72). It follows that there exist constants $C, \sigma > 0$, which do not depend on φ , such that

$$\forall \boldsymbol{u} \in e^{-i\boldsymbol{\varphi}} \boldsymbol{\Omega}_{-\psi} \cap \tilde{\boldsymbol{B}}, \ \forall \boldsymbol{\lambda} \in \boldsymbol{R}_{\boldsymbol{\varphi}}, \quad \left| e^{-\boldsymbol{\lambda} \cdot \boldsymbol{u}} I(\boldsymbol{\lambda}) \right| \le C^{|u_1| + \dots + |u_d|} \cdot e^{-\sigma(|\lambda_1 u_1| + \dots + |\lambda_d u_d|)}. \tag{73}$$

One can easily check by direct computation that

$$J := \sup_{\boldsymbol{\varphi} \in (-\delta, \delta)^d} \sup_{\boldsymbol{u} \in \tilde{\boldsymbol{B}}} \int_{\boldsymbol{R}_{\boldsymbol{\varphi}}} e^{-\sigma(|\lambda_1 u_1| + \dots + |\lambda_d u_d|)} |\mathrm{d}\boldsymbol{\lambda}|$$
(74)

is finite. It follows that $|\mathcal{L}_{\boldsymbol{c}}[I](\boldsymbol{u})| \leq J \cdot C^{|u_1|+\dots+|u_d|}$ for all $\boldsymbol{u} \in e^{-i\varphi} \Omega_{-\psi} \cap \tilde{\boldsymbol{B}}$. This implies (70), because J, C are independent of φ , and we have $\Omega_{\delta''} = \bigcup_{\varphi \in (-\delta, \delta)^d} e^{-i\varphi} \Omega_{-\psi}$ with the choice $\psi := \delta - \delta''$.

Now assume $\delta^{\circ} < \delta' < \delta''$ and consider the contour $V_{\delta'} = (\partial U_{\delta',1})^d$, where $U_{\delta',r}$ is defined above (55) (see also Figure 2(a)). Then (55) implies that

$$\forall \boldsymbol{\lambda} \in \boldsymbol{K}_{\delta^{\circ}}, \; \forall \boldsymbol{u} \in \boldsymbol{V}_{\delta'}, \quad \left| e^{\boldsymbol{\lambda} \cdot \boldsymbol{u}} \right| \le C_*^{|\lambda_1| + \dots + |\lambda_d|} e^{-\sigma_*(|\lambda_1 u_1| + \dots + |\lambda_d u_d|)} \tag{75}$$

Since we have $V_{\delta'} \subset \Omega_{\delta''} \cap \dot{B}$, the last display and (70) imply that $\mathcal{B}[\mathcal{L}_{c}[I]]$ is well-defined and analytic on $\{\lambda \in K_{\delta^{\circ}} : \forall j, |\lambda_{j}| > \frac{m+1}{\sigma_{*}}\}$.

Proof of Theorem 2(3). Notice that the Borel transform \mathcal{B} and the Laplace transform \mathcal{L}_{c} can both be written as the product of d univariate integral transforms: we have $\mathcal{B} = \mathcal{B}^{(1)} \circ \cdots \circ \mathcal{B}^{(d)}$ and $\mathcal{L}_{c} = \mathcal{L}_{c_{1}}^{(1)} \circ \cdots \circ \mathcal{L}_{c_{d}}^{(1)}$ with

$$\mathcal{B}^{(j)}[H](\boldsymbol{u}_{\hat{j}}) := \frac{1}{2\pi i} \int_{V_{\delta'}} e^{\lambda u_j} H(\boldsymbol{u}) \mathrm{d}u_j \quad \text{and} \quad \mathcal{L}_c^{(j)}[I](\boldsymbol{\lambda}_{\hat{j}}) := \int_c^\infty e^{\lambda_j u} I(\boldsymbol{\lambda}) \mathrm{d}\lambda_j, \tag{76}$$

where $\boldsymbol{u}_{\hat{j}} = (u_1, \ldots, u_{j-1}, \lambda, u_{j+1}, \ldots, u_d)$ and $\boldsymbol{\lambda}_{\hat{j}} = (\lambda_1, \ldots, \lambda_{j-1}, u, \lambda_{j+1}, \ldots, \lambda_d)$. Moreover, for all $j \neq k$, $\mathcal{B}^{(j)}$ commute with $\mathcal{L}_{c_k}^{(k)}$, because they operate on different variables. It follows that

$$\mathcal{B} \circ \mathcal{L}_{\boldsymbol{c}} = \left(\mathcal{B}^{(1)} \circ \mathcal{L}^{(1)}_{c_1} \right) \circ \cdots \circ \left(\mathcal{B}^{(d)} \circ \mathcal{L}^{(d)}_{c_d} \right) \quad \text{and} \quad \mathcal{L}_{\boldsymbol{c}} \circ \mathcal{B} = \left(\mathcal{L}^{(1)}_{c_1} \circ \mathcal{B}^{(1)} \right) \circ \cdots \circ \left(\mathcal{L}^{(d)}_{c_d} \circ \mathcal{B}^{(d)} \right).$$
(77)

The above decompositions allow us to reduce the proof of Theorem 2(3) to the univariate case. In the following, we assume d = 1 and fix some $\delta \in (0, \frac{\pi}{2})$ and $c \in \mathbb{R}_{>0}$.

Let $I \in \mathcal{P}(K_{\delta})$. Let us show that $\mathcal{B} \circ \mathcal{L}_{c}[I](\lambda) = I(\lambda)$ for all $\lambda > c$. For this, fix some $\delta' \in (0, \delta)$ and let $V_{\delta'} = \partial \Omega_{\delta}$ with the domain Ω_{δ} defined above (55). Let $V_{\delta'}^{+}$ (resp. $V_{\delta'}^{-}$) be the part of $V_{\delta'}$ in the upper half-plane (resp. lower half-plane). Recall from the proof of Theorem 2(2) that $\mathcal{L}_{c}[I]$ can be expressed as an integral over $R_{c,\varphi}$ for any $\varphi \in (-\delta, \delta)$, where $R_{c,\varphi}$ is the contour shown in Figure 2(3). Fix some $\varphi \in (\delta', \delta)$. Then we can write

$$\mathcal{B} \circ \mathcal{L}_{c}[I](\lambda) = \frac{1}{2\pi i} \int_{V_{\delta'}} e^{\lambda u} \mathcal{L}_{c}[I](u) du$$

= $\frac{1}{2\pi i} \left(\int_{V_{\delta'}^{+}} e^{\lambda u} \left(\int_{R_{c,-\varphi}} e^{-\tau u} I(\tau) d\tau \right) du + \int_{V_{\delta'}^{-}} e^{\lambda u} \left(\int_{R_{c,\varphi}} e^{-\tau u} I(\tau) d\tau \right) du \right).$

Thanks to the bounds (55) and (67), it is not hard to see that $(u, \tau) \mapsto e^{\lambda u} e^{-\tau u} I(\tau)$ is integrable on both $V_{\delta'}^+ \times R_{c,-\varphi}$ and $V_{\delta'}^- \times R_{c,\varphi}$. Thus we have by Fubini's theorem

$$\begin{aligned} \mathcal{B} \circ \mathcal{L}_{c}[I](\lambda) &= \frac{1}{2\pi i} \left(\int_{R_{c,-\varphi}} I(\tau) \left(\int_{V_{\delta'}^{+}} e^{\lambda u - \tau u} \mathrm{d}u \right) \mathrm{d}\tau + \int_{R_{c,\varphi}} I(\tau) \left(\int_{V_{\delta'}^{-}} e^{\lambda u - \tau u} \mathrm{d}u \right) \mathrm{d}\tau \right) \\ &= \frac{1}{2\pi i} \left(\int_{R_{c,-\varphi}} I(\tau) \cdot \left[\frac{e^{(\lambda - \tau)u}}{\lambda - \tau} \right]_{1}^{e^{i(\frac{\pi}{2} + \delta')\infty}} \mathrm{d}\tau + \int_{R_{c,\varphi}} I(\tau) \cdot \left[\frac{e^{(\lambda - \tau)u}}{\lambda - \tau} \right]_{e^{-i(\frac{\pi}{2} + \delta')\infty}}^{1} \mathrm{d}\tau \right) \\ &= \frac{1}{2\pi i} \left(\int_{R_{c,-\varphi}} I(\tau) \cdot \frac{e^{\lambda - \tau}}{\tau - \lambda} \mathrm{d}\tau + \int_{R_{c,\varphi}} I(\tau) \cdot \frac{e^{\lambda - \tau}}{\lambda - \tau} \mathrm{d}\tau \right) \\ &= \frac{1}{2\pi i} \int_{\mathcal{R}} I(\tau) \cdot \frac{e^{\lambda - \tau}}{\tau - \lambda} \mathrm{d}\tau \end{aligned}$$

where $\mathcal{R} = R_{c,-\varphi} \cup \tilde{R}_{c,\varphi}$ and $\tilde{R}_{c,\varphi}$ is the contour $R_{c,\varphi}$ oriented in the opposite direction. Notice that $\mathcal{R} \subset K_{\delta}$ is a bi-infinite contour whose two ends extend to ∞ . Moreover, the integrand $\tau \mapsto I(\tau) \frac{e^{\lambda-\tau}}{\tau-\lambda}$ is meromorphic on K_{δ} , has a unique simple pole at $\tau = \lambda$ (which is on the right of the contour \mathcal{R}), and decays exponentially when $\Re \mathfrak{e}(\tau) \to \infty$. It follows from the residue theorem that

$$\frac{1}{2\pi i} \int_{\mathcal{R}} I(\tau) \cdot \frac{e^{\lambda - \tau}}{\tau - \lambda} d\tau = \operatorname{Res}_{\tau \to \lambda} \left(I(\tau) \frac{e^{\lambda - \tau}}{\tau - \lambda} \right) = I(\lambda) \,. \tag{78}$$

This proves $\mathcal{B} \circ \mathcal{L}_c[I](\lambda) = I(\lambda)$ for all $\lambda > c$, and thus for all $\lambda \in K_{\delta}$ by analytic continuation. In other words, $\mathcal{B} \circ \mathcal{L}_c$ is the identity on $\mathcal{P}(K_{\delta})$. By the decomposition (77), the same is true in any dimension d.

Now fix some $H \in \mathcal{P}(\Omega_{\delta})$ and u > 0. Consider the contour $V_{\delta'} = \partial U_{\delta',r}$ for some $\delta' \in (0, \delta)$ and $r \in (0, u)$. It is not hard to check that the function $(\lambda, v) \mapsto e^{-\lambda u} e^{\lambda v} H(v)$ is integrable on $[c, \infty) \times V_{\delta'}$. By Fubini's theorem, we have

$$\mathcal{L}_c \circ \mathcal{B}[H](u) = \frac{1}{2\pi i} \int_c^\infty e^{-\lambda u} \left(\int_{V_{\delta'}} e^{\lambda v} H(v) \mathrm{d}v \right) \mathrm{d}\lambda$$
(79)

$$= \frac{1}{2\pi i} \int_{V_{\delta'}} H(v) \left(\int_c^\infty e^{-(u-v)\lambda} \mathrm{d}\lambda \right) \mathrm{d}v = \frac{1}{2\pi i} \int_{V_{\delta'}} H(v) \frac{e^{-(u-v)c}}{u-v} \mathrm{d}v \,. \tag{80}$$

For each $u \in \mathbb{C}$, let

$$E_c(u) = -\frac{1}{2\pi i} \lim_{R \to \infty} \int_{\partial U_{\delta',R}} H(v) \frac{e^{-(u-v)c}}{u-v} \mathrm{d}v \,. \tag{81}$$

It is not hard to see that the above limit stablizes when R > |u|, and defines an entire function of u. Actually, by the residue theorem, we have for all R > |u|,

$$\frac{1}{2\pi i} \int_{\partial U_{\delta',R}} H(v) \frac{e^{-(u-v)c}}{u-v} dv = \frac{1}{2\pi i} \int_{V_{\delta'}} H(v) \frac{e^{-(u-v)c}}{u-v} dv + \operatorname{Res}_{v \to u} H(v) \frac{e^{-(u-v)c}}{u-v} dv + \operatorname$$

It follows that $\mathcal{L}_c \circ \mathcal{B}[H](u) = H(u) + E_c(u)$ for all u > 0. By analytic continuation, the same is true for all $u \in \Omega_{\delta}$. To recover the case of general dimension d, we apply this formula of $\mathcal{L}_c \circ \mathcal{B}[H]$ to each factor in the decomposition (77). To simplify notation, we write $\mathcal{I}_j = \mathcal{L}_{c_j}^{(j)} \circ \mathcal{B}^{(j)}$, then

$$\mathcal{L}_{\boldsymbol{c}} \circ \mathcal{B}[H] = (\mathcal{I}_{1} \circ \cdots \circ \mathcal{I}_{d-1} \circ \mathcal{I}_{d}) [H]$$

= $(\mathcal{I}_{1} \circ \cdots \circ \mathcal{I}_{d-1}) [H] + (\mathcal{I}_{1} \circ \cdots \circ \mathcal{I}_{d-1}) [E_{\boldsymbol{c}}^{(d)}]$
= $\cdots \qquad \cdots$
= $H + E_{\boldsymbol{c}}^{(1)} + \mathcal{I}_{1}[E_{\boldsymbol{c}}^{(2)}] + \cdots + (\mathcal{I}_{1} \circ \cdots \circ \mathcal{I}_{d-1}) [E_{\boldsymbol{c}}^{(d)}].$

where each $E_{\boldsymbol{c}}^{(j)}$ is an analytic function on $\boldsymbol{\Omega}_{\delta}$ which is entire with respect to the variable u_j . The same is true for $\mathcal{I}_1 \circ \cdots \mathcal{I}_{j-1}[E_{\boldsymbol{c}}^{(j)}]$, because the integral transform $\mathcal{I}_1 \circ \cdots \mathcal{I}_{j-1}$ does not affect the variable u_j . It follows that $E_{\boldsymbol{c}}^{(1)} + \mathcal{I}_1[E_{\boldsymbol{c}}^{(2)}] + \cdots + (\mathcal{I}_1 \circ \cdots \circ \mathcal{I}_{d-1})[E_{\boldsymbol{c}}^{(d)}]$ is a demi-entire function on $\boldsymbol{\Omega}_{\delta}$. This proves the decomposition $\mathcal{L}_{\boldsymbol{c}} \circ \mathcal{B}[H] = H + E_{\boldsymbol{c}}$ in any dimension d, and concludes the proof of Theorem 2.

4 Discussions

In this section, we discuss some additional properties of the classes of multivariate functions used in the statement of the multivariate transfer theorems.

 Δ -analytic functions. Like in the univariate case, the multivariate transfer theorems relies fundamentally on the Δ -analyticity of the generating function. But unlike the univariate case, Δ -analyticity is a rather restrictive condition in the multivariate setting. In particular, it implies that the dominant singularity (for any reasonable definition of the term) is unique and independent of the exponent θ and the direction λ of the θ -diagonal limit taken. This is in stark contrast with the case of rational functions, where the dominant singularities (a.k.a. contributing critical points) generically depend on the direction of the diagonal limit taken. In fact, a multivariate rational function is never Δ -analytic at any of its poles, unless its denominator has a univariate linear factor which vanishes at this pole. **Proposition 5** (Genuinely multivariate rational functions are never Δ -analytic). If a rational function is Δ -analytic at **1** and has a pole at **1**, then its denominator is divisible by $z_j - 1$ for some $1 \le j \le d$.

Remark. The proposition as well as its proof given below can be easily generalized to meromorphic functions. We shall not enter into the details here and refer to [8, Section 3.1.1] for an introduction to meromorphic functions in several variable.

Proof. Consider a rational function $A = \frac{F}{G}$, where F, G are coprime polynomials. The proposition follows directly from Lemma 6 below about the local geometry of the zero set of a polynomial: it suffices to apply Lemma 6 to the polynomial $H(\boldsymbol{u}) = G(1 - \boldsymbol{u})$.

For
$$\delta, \epsilon > 0$$
, let $\Omega_{\delta,\epsilon} = (\Omega_{\delta,\epsilon})^d$ with $\Omega_{\delta,\epsilon} := \Omega_{\delta} \cap B_{0,\epsilon} \equiv \left\{ u \in \mathbb{C} : 0 < |u| < \epsilon \text{ and } |\arg(u)| < \frac{\pi}{2} + \delta \right\}$.

Lemma 6 (Local version of Proposition 5). If a polynomial $H \in \mathbb{C}[\mathbf{u}]$ has a zero at $\mathbf{0}$, but no zero on $\Omega_{\delta,\epsilon}$ for some $\delta, \epsilon > 0$, then $H(\mathbf{u})$ is divisible by u_j for some $1 \leq j \leq d$.

Proof. When d = 1, the lemma is trivial: a univariate polynomial H(u) has a zero at 0 if and only if it is divisible by u.

When d = 2, consider a polynomial $H \in \mathbb{C}[u, v]$ not divisible by u nor v, such that H(0, 0) = 0. According to the Newton-Puiseux theorem (see e.g. [2, Corollary 1.5.5, Theorem 1.7.2]), there exist an integer $n \geq 1$ and a nonzero analytic function s defined on a neighborhood of 0, such that for every determination of the n-th root, the mapping $\varphi(u) = s(u^{1/n})$ satisfies $H(u, \varphi(u)) = 0$ for all u in some neighborhood of 0. Up to swapping the variables u and v, we can assume without loss of generality that $\varphi(u) = O(u)$ as $u \to 0$. Since s is analytic at 0 and is not identically zero, there exist $c \neq 0$ and $\alpha \geq 1$ such that $\varphi(u) \sim c \cdot u^{\alpha}$ as $u \to 0$. Geometrically, this means that the mapping $\varphi(\Omega_{\delta,\epsilon})$ conains an angle of $\alpha(\pi + 2\delta)$ at 0. Since $\alpha \geq 1$, we have $\alpha(\pi + 2\delta) + (\pi + 2\delta) > 2\pi$. This implies that $\varphi(\Omega_{\delta,\epsilon}) \cap \Omega_{\delta,\epsilon} \neq \emptyset$. In other words, the graph of the mapping φ intersects $(\Omega_{\delta,\epsilon})^2$. It follows that the polynomial H has zeros on $(\Omega_{\delta,\epsilon})^2$ for any $\delta, \epsilon > 0$. By contraposition, this proves Lemma 6 when d = 2.

For $d \ge 3$, we give a proof by contradiction based on the result of the case d = 2: Let $H \in \mathbb{C}[\boldsymbol{u}]$ be a polynomial with a zero at $\boldsymbol{0}$, no zero on $\boldsymbol{\Omega}_{\delta,\epsilon}$ for some $\delta, \epsilon > 0$, and such that $H(\boldsymbol{u})$ is not divisible by u_j for any $1 \le j \le d$. For $1 \le j \le d$ and $m \ge 0$, define $E_m^{(j)} = \{\boldsymbol{n} \in \mathbb{N}^d \mid (n_1 + \cdots + n_d) - n_j = m\}$. Let us prove the following statement by induction on m:

$$\forall j \in \{1, \dots, d\}, \, \forall \boldsymbol{n} \in E_m^{(j)}, \quad [\boldsymbol{u}^n] H(\boldsymbol{u}) = 0.$$
 (\mathcal{H}_m)

For $\boldsymbol{x} \in \mathbb{R}^{d-1}_{>0}$, let $H_{\boldsymbol{x}}(u, v) = H(u, vx_1, \dots, vx_{d-1})$. One can check that for all $k \in \mathbb{N}$, we have

$$[v^k]H_{\boldsymbol{x}}(u,v) = \sum_{\boldsymbol{n}\in E_k^{(1)}} [\boldsymbol{u}^n]H(\boldsymbol{u}) \cdot u^{n_1} \cdot x_1^{n_2} \cdots x_{d-1}^{n_d}.$$
(82)

Now fix some $m \ge 0$ and assume that (\mathcal{H}_k) is true for all $0 \le k < m$. By (82), we have $[v^k]H_{\boldsymbol{x}}(u,v) = 0$ for all k < m. Hence $H_{\boldsymbol{x}}^{(m)}(u,v) := v^{-m}H_{\boldsymbol{x}}(u,v)$ is a polynomial in (u,v). Moreover, $H_{\boldsymbol{x}}^{(m)}(0,0) = 0$: For m = 0, we have $H_{\boldsymbol{x}}^{(0)}(0,0) = H_{\boldsymbol{x}}(0,0) = H(\mathbf{0}) = 0$ by assumption. When $m \ge 1$, one can check that

$$\left\{ \boldsymbol{n} \in E_m^{(1)} \, \big| \, n_1 = 0 \right\} \subseteq \bigcup_{k < m} \left(\bigcup_{j=1}^d E_k^{(j)} \right).$$
(83)

Hence the hypotheses $(\mathcal{H}_k)_{k < m}$ and (82) imply that $H_{\boldsymbol{x}}^{(m)}(0,0) \equiv [u^0 v^m] H_{\boldsymbol{x}}(u,v) = 0$ as well. On the other hand, since $H(\boldsymbol{u}) \neq 0$ for all $\boldsymbol{u} \in \Omega_{\delta,\epsilon}$, we have $H_{\boldsymbol{x}}^{(m)}(u,v) = v^{-m}H(u,vx_1,\ldots,vx_{d-1}) \neq 0$ for all $(u,v) \in (\Omega_{\delta,\tilde{\epsilon}})^2$, where $\tilde{\epsilon} = \min(1, x_1^{-1}, \ldots, x_{d-1}^{-1}) \cdot \epsilon > 0$. So, according to the result of the case d = 2, the polynomial $H_{\boldsymbol{x}}^{(m)}(u,v)$ is either divisible by u or divisible by v. In particular, we have

$$H_{\boldsymbol{x}}^{(m)}(0,1) \cdot H_{\boldsymbol{x}}^{(m)}(u,0) \equiv H(0,\boldsymbol{x}) \cdot [v^m] H_{\boldsymbol{x}}(u,v) = 0$$
(84)

for all $\boldsymbol{x} \in \mathbb{R}_{>0}^{d-1}$. By analytic continuation, the above identity is valid for all $\boldsymbol{x} \in \mathbb{C}^{d-1}$. Because $H(\boldsymbol{u})$ is not divisible by u_1 , the polynomial $\boldsymbol{x} \mapsto H(0, \boldsymbol{x})$ is not identically zero. Since the ring of polynomials is an integral domain, this implies that $[\boldsymbol{v}^m]H_{\boldsymbol{x}}(u, v) = 0$ as a polynomial in \boldsymbol{x} . Comparing this to (82), we see that $[\boldsymbol{u}^n]H(\boldsymbol{u}) = 0$ for all $\boldsymbol{n} \in E_m^{(1)}$. The same argument works for $E_m^{(j)}$ for any $1 \leq j \leq d$. Therefore (\mathcal{H}_m) is true.

By induction, (\mathcal{H}_m) holds for all $m \geq 0$. But this implies that $H(\mathbf{u}) \equiv 0$, which contradicts the assumption that $H(\mathbf{u})$ is not divisible by u_j . This completes the proof of Lemma 6 for general $d \geq 3$. \Box

Demi-analytic functions. A univariate function A is demi-analytic at ρ if and only if it is analytic in a neighborhood of ρ , and it is demi-entire if and only if it is an entire function. Thus the names "demi-analytic" and "demi-entire".

It is easy to check that a (θ_0, θ) -homogeneous function is demi-entire with respect to a cone if and only if it is demi-analytic at **0** with respect to the same cone.

Given the form of the asymptotic expansion (17) and the analyticity of the scaling function $I(\lambda)$, it is not hard to see that the coefficients $[\mathbf{z}^n]A(\mathbf{z})$ decays exponentially as $n_0 \to \infty$ if and only if Iis identically zero. By Corollary 3, this happens if and only if the homogeneous component $H(\mathbf{u})$ is demi-analytic. In this sense, the demi-analyticity condition in Theorem 1(1) is optimal.

Generalized homogeneous functions. If H is a (θ_0, θ) -homogeneous function, then for all $u_1 > 0$, we have

$$H(\boldsymbol{u}) = u_1^{\theta_0/\theta_1} \cdot H\left(1, \frac{u_2}{u_1^{\theta_2/\theta_1}}, \cdots, \frac{u_d}{u_1^{\theta_d/\theta_1}}\right).$$
(85)

Inversely, for any function $h(r_2, \ldots, r_d)$ of d-1 variables, $H(\boldsymbol{u}) = u^{\theta_0/\theta_1} \cdot h(u_1^{-\theta_2/\theta_1}u_2, \cdots, u_1^{-\theta_d/\theta_1}u_d)$ is a $(\theta_0, \boldsymbol{\theta})$ -homogeneous function satisfying $h(r_2, \ldots, r_d) = H(1, r_2, \ldots, r_d)$. In particular, when d = 1, the only (θ_0, θ_1) -homogeneous analytic functions are the power functions $H(\boldsymbol{u}) = C u^{\theta_0/\theta_1}$ ($C \in \mathbb{C}$).

For each $\boldsymbol{\alpha} \in \mathbb{R}^d$, the monomial $\boldsymbol{u} \mapsto \boldsymbol{u}^{\boldsymbol{\alpha}}$ is $(\theta_0, \boldsymbol{\theta})$ -homogeneous for all $(\theta_0, \boldsymbol{\theta})$ such that $\boldsymbol{\alpha} \cdot \boldsymbol{\theta} = \theta_0$. Clearly, a linear combination of finitely many such monomials is also $(\theta_0, \boldsymbol{\theta})$ -homogeneous, as long as their multi-exponents $\boldsymbol{\alpha}$ satisfy the above linear relation for the same $(\theta_0, \boldsymbol{\theta})$. Such polynomials are obviously *not* all the homogeneous functions, but they provide a good intuition for how a non-homogeneous function could be decomposed into homogeneous components.

Functions of polynomial type. In Theorem 1', the remainder term in the asymptotic expansion of the generating function A(z) was expressed as a big-O of some (θ_0, θ) -homogeneous function of polynomial type. The following lemma tells us that every such remainder term can also be bounded by finitely many *monomials* with the same homogeneity.

Lemma 7. Let K be a cone in \mathbb{C}^d . If $H : K \to \mathbb{C}$ is a (θ_0, θ) -homogeneous function of polynomial type locally at $\mathbf{0}$, then there exist a constant C and finitely many vectors $\boldsymbol{\alpha}^{(k)} \in \mathbb{R}^d$, such that $\boldsymbol{\alpha}^{(k)} \cdot \boldsymbol{\theta} = \theta_0$ for all k, and

$$\forall \boldsymbol{u} \in \boldsymbol{K}, \quad |H(\boldsymbol{u})| \leq C \cdot \sum_{k} \left| \boldsymbol{u}^{\boldsymbol{\alpha}^{(k)}} \right|.$$
(86)

Proof. By considering $\mathcal{H}(s) = H(s^{\theta})$ instead of H, we can assume without loss of generality that $\theta = 1$. By assumption, there exist constants $C, M, \epsilon > 0$ such that on $\{u \in \mathbf{K} | \forall 1 \leq j \leq d, |u_j| < \epsilon\}$, we have

$$|H(\boldsymbol{u})| \le C \cdot \left(|u_1|^{-M} + \dots + |u_d|^{-M}\right)$$
 (87)

Let i be such that $|u_i| = \max_{1 \le j \le d} |u_j|$. Then we can rescale the vector \boldsymbol{u} by $\frac{\epsilon}{|u_i|}$ to place it in the above

set. Since *H* is $(\theta_0, \mathbf{1})$ -homogeneous, we have $H(\mathbf{u}) = \epsilon^{-\theta_0} |u_i|^{\theta_0} H\left(\epsilon \frac{u_1}{|u_i|}, \ldots, \epsilon \frac{u_d}{|u_i|}\right)$, and therefore

$$|H(\boldsymbol{u})| \leq \epsilon^{-\theta_0} |u_i|^{\theta_0} C \cdot \left(\epsilon^{-M} \left|\frac{u_1}{u_i}\right|^{-M} + \dots + \epsilon^{-M} \left|\frac{u_d}{u_i}\right|^{-M}\right) \\ = C \epsilon^{-\theta_0 - M} \left(|u_i^{\theta_0 + M} u_2^{-M}| + \dots + |u_i^{\theta_0 + M} u_d^{-M}| \right)$$

Summing the right hand side over i gives a bound of H(u) of the form (86) for all $u \in K$.

Lemma 7 was briefly used in the proof of Theorem 1' to obtain the bound (26). It implies (26) because the variable $u \in V$ in (26) satisfies $|u_j| \ge 1$ for all $1 \le j \le d$.

In practical examples, it is usually not hard to check that a function is of polynomial type, in particular thanks to the following lemma.

Lemma 8 (Closure properties of $\mathcal{P}(\mathbf{K}_{\delta})$ and $\mathcal{P}(\mathbf{\Delta}_{\delta})$). For any $\delta \in (0, \pi)$, the space $\mathcal{P}(\mathbf{K}_{\delta})$ forms a \mathbb{C} -algebra with respect to pointwise addition and multiplication of functions. Moreover, if f is any function such that |f(x)| is bounded by a polynomial of |x| (e.g. $f(x) = x^{\beta}, \beta \in \mathbb{C}$),

then for all $h \in \mathcal{P}(\mathbf{K}_{\delta})$ such that $f \circ h$ is well-defined and analytic on \mathbf{K}_{δ} , we have $f \circ h \in \mathcal{P}(\mathbf{K}_{\delta})$. The same is true for $\mathcal{P}(\mathbf{\Delta}_{\delta})$ with $\delta \in (0, \pi/2)$.

Proof. Fix $\delta \in (0, \pi)$ and let $g, h \in \mathcal{P}(\mathbf{K}_{\delta})$. By definition there exist C, M > 0 such that

$$|g(\boldsymbol{u})| \le C \cdot \left(|u_1|^M + |u_1|^{-M} + \dots + |u_d|^M + |u_d|^{-M}\right)$$
(88)

and
$$|h(\boldsymbol{u})| \le C \cdot \left(|u_1|^M + |u_1|^{-M} + \dots + |u_d|^M + |u_d|^{-M}\right)$$
 (89)

for all $\boldsymbol{u} \in \boldsymbol{K}_{\delta}$. It is clear that any linear combination of g and h satisfies a bound of the same form. So $\mathcal{P}(\boldsymbol{K}_{\delta})$ is a \mathbb{C} -vector space. In addition, since $ab \leq \frac{1}{2}(a^2 + b^2)$ for all $a, b \in \mathbb{R}$, we have

$$\begin{aligned} |g(\boldsymbol{u})h(\boldsymbol{u})| &\leq C^2 \sum_{i,j=1}^d \left(|u_i|^M |u_j|^M + |u_i|^M |u_j|^{-M} + |u_i|^{-M} |u_j|^M + |u_i|^{-M} |u_j|^{-M} \right) \\ &\leq C^2 \sum_{i,j=1}^d \left(|u_i|^{2M} + |u_j|^{2M} + |u_i|^{-2M} + |u_j|^{-2M} \right). \end{aligned}$$

So $\mathcal{P}(\mathbf{K}_{\delta})$ is also closed under multiplication, and therefore forms a \mathbb{C} -algebra.

Now let f be a function such that |f(x)| is bounded by a polynomial of |x|. Then there exist c > 0and $m \in \mathbb{Z}_{\geq 0}$ such that $|f(x)| \leq c \cdot (1 + |x|^m)$. If $h \in \mathcal{P}(\mathbf{K}_{\delta})$ and $f \circ h$ is well-defined and analytic on \mathbf{K}_{δ} , then we have $h^m \in \mathcal{P}(\mathbf{K}_{\delta})$ by the closure of $\mathcal{P}(\mathbf{K}_{\delta})$ under multiplication, and therefore

$$|f \circ h(\boldsymbol{u})| \le c \cdot (1 + |h(\boldsymbol{u})^{m}|) \le c \cdot \left(1 + C \cdot \left(|u_{1}|^{M} + |u_{1}|^{-M} + \dots + |u_{d}|^{M} + |u_{d}|^{-M}\right)\right)$$
(90)

for some C, M > 0. It follows that $f \circ h \in \mathcal{P}(\mathbf{K}_{\delta})$. The same proof works for $\mathcal{P}(\mathbf{\Delta}_{\delta})$.

Let $\Omega \subseteq \mathbb{C}^d$. A function $A : \Omega \to \mathbb{C}$ is algebraic if there exists a polynomial $E \in \mathbb{C}[a, \mathbf{z}]$, such that $E(A(\mathbf{z}), \mathbf{z}) = 0$ for all $\mathbf{z} \in \Omega$.

Lemma 9 (Algebraic functions are of polynomial type). If A is an algebraic function analytic on Δ_{δ} , then $A \in \mathcal{P}(\Delta_{\delta'})$ for all $\delta' \in (0, \delta)$.

Proof. Let A be an algebraic function which is analytic on Δ_{δ} . For $z \in \Delta_{\delta}$, let

$$h(\boldsymbol{z}) = \frac{1}{\operatorname{dist}(\boldsymbol{z}, \mathbb{C}^d \setminus \boldsymbol{\Delta}_{\delta})} \equiv \left(\inf_{\boldsymbol{w} \in \mathbb{C}^d \setminus \boldsymbol{\Delta}_{\delta}} \|\boldsymbol{z} - \boldsymbol{w}\|\right)^{-1},\tag{91}$$

where $\|\boldsymbol{z} - \boldsymbol{w}\|$ is the Euclidean distance between \boldsymbol{z} and \boldsymbol{w} , viewed as points in \mathbb{R}^{2d} . Let us show that there exist C, M > 0 such that

$$\forall \boldsymbol{z} \in \boldsymbol{\Delta}_{\delta}, \quad |A(\boldsymbol{z})| \le C \cdot h(\boldsymbol{z})^{M}.$$
(92)

This is a variant of the *Lojasiewicz inequality* in semi-algebraic geometry. See e.g. [1, Chapter 2] for an introduction. Recall that a set $S \subseteq \mathbb{R}^n$ is semi-algebraic if and only if it can be defined by a first order logic formula involving only polynomial conditions (i.e. equations or inequalities) and quantifiers \forall/\exists over \mathbb{R} . And a function $f: S \to \mathbb{R}^m$ is semi-algebraic if and only if $\operatorname{Graph}(f) := \{(\boldsymbol{x}, f(\boldsymbol{x})) \in \mathbb{R}^{n+m} | \boldsymbol{x} \in S\}$ is a semi-algebraic set.

It is not hard to write a first order formula that describes the set Δ_{δ} (viewed as a subset of \mathbb{R}^{2d}). Therefore Δ_{δ} is a semi-algebraic set. It is well-known that a continuous algebraic function defined on an open semi-algebraic set is always semi-algebraic (see e.g. [10, Theorem 11]). Hence the function $A : \Delta_{\delta} \to \mathbb{C} \equiv \mathbb{R}^2$ is semi-algebraic. Thanks to general closure properties of the class of semi-algebraic functions (c.f. [1, Section 2.2]), we deduce from the above facts that the functions $h : \Delta_{\delta} \to \mathbb{R}$ and $|A| : \Delta_{\delta} \to \mathbb{R}$ are also semi-algebraic.

For each t > 0, let $G_t = \{ z \in \Delta_{\delta} : h(z) = t \}$ and $g(t) = \sup_{z \in G_t} |A(z)|$, with the convention that $\sup \emptyset = 0$. The graph of the function $g : \mathbb{R}_{>0} \to \mathbb{R}$ can be described by a first order formula as follows:

$$\begin{aligned} \operatorname{Graph}(g) &= \left\{ \left(t, y \right) \in \mathbb{R}^2 \ \Big| \ t > 0 \text{ and } \left(\forall \boldsymbol{z} \in G_t, \ |A(\boldsymbol{z})| \le y \right) \\ & \text{and } \left(\left(G_t = \varnothing \text{ and } y = 0 \right) \text{ or } \left(\exists \boldsymbol{z} \in G_t, \ |A(\boldsymbol{z})| = y \right) \right) \right\} \end{aligned}$$

where

$$\begin{split} \left(\forall \boldsymbol{z} \in G_t, \, |A(\boldsymbol{z})| \leq y \right) \ \Leftrightarrow \ \left(\forall \boldsymbol{z} \in \mathbb{R}^{2d}, \left(\boldsymbol{z} \in \boldsymbol{\Delta}_{\delta} \text{ and } (\boldsymbol{z}, t) \in \mathtt{Graph}(h) \right) \\ \Rightarrow \left(\exists w \in \mathbb{R}, \, (\boldsymbol{z}, w) \in \mathtt{Graph}(|A|) \text{ and } w \leq y \right) \right), \\ \left(G_t = \varnothing \text{ and } y = 0 \right) \Leftrightarrow \ \left(\forall \boldsymbol{z} \in \mathbb{R}^{2d}, \left(\boldsymbol{z} \in \boldsymbol{\Delta}_{\delta} \Rightarrow (\boldsymbol{z}, t) \notin \mathtt{Graph}(h) \right) \right), \\ \left(\exists \boldsymbol{z} \in G_t, \, |A(\boldsymbol{z})| = y \right) \Leftrightarrow \ \left(\exists \boldsymbol{z} \in \mathbb{R}^{2d}, \left(\boldsymbol{z} \in \boldsymbol{\Delta}_{\delta} \text{ and } (\boldsymbol{z}, t) \in \mathtt{Graph}(h) \text{ and } (\boldsymbol{z}, y) \in \mathtt{Graph}(|A|) \right) \right). \end{split}$$

Since h and |A| are semi-algebraic functions, the conditions $(z, t) \in \operatorname{Graph}(h)$ and $(z, y) \in \operatorname{Graph}(|A|)$ can be further expanded into first order logic formulas involving only polynomial conditions and quantifiers over \mathbb{R} . It follows that g is a semi-algebraic function. A classical result on the growth rate of univariate semi-algebraic functions [1, Proposition 2.6.1] states that there exist constants $C, M, t_0 > 0$ such that $g(t) \leq C \cdot t^M$ for all $t > t_0$. The definition of g implies that it is bounded on $(0, t_0]$. Hence the bound $g(t) \leq C \cdot t^M$ extends to all t > 0. This proves (92).

Finally, notice that for all $\delta' \in (0, \delta)$, there exists a constant c > 0 such that

$$\forall \boldsymbol{z} \in \boldsymbol{\Delta}_{\delta'}, \quad c^{-1} \cdot \operatorname{dist}(\boldsymbol{z}, \mathbb{C}^d \setminus \boldsymbol{\Delta}_{\delta}) \leq \min_{1 \leq j \leq d} |z_j - 1| \leq c \cdot \operatorname{dist}(\boldsymbol{z}, \mathbb{C}^d \setminus \boldsymbol{\Delta}_{\delta}).$$
(93)

Moreover, we have

$$\left(\min_{1 \le j \le d} |z_j|\right)^{-M} \le |z_1 - 1|^{-M} + \dots + |z_d - 1|^{-M} \le d \cdot \left(\min_{1 \le j \le d} |z_j|\right)^{-M}.$$
(94)

It follows that $|A(\mathbf{z})| \leq \tilde{C} \sum_{j=1}^{d} |z_j - 1|^{-M}$ on $\mathbf{\Delta}_{\delta'}$ for some $\tilde{C}, M > 0$. That is, $A \in \mathcal{P}(\mathbf{\Delta}_{\delta'})$.

Background on ACSV and its general strategy. The results in paper fall under the topic of *analytic combinatorics in several variables* (ACSV). The remaining paragraphs provide some background on ACSV and where this work stands relative to the others.

In general, analytic combinatorics aims at understanding the enumerative properties of large combinatorial structures through the analytic properties of their generating functions. This is usually done in two steps: First, some (possibly implicit) expression of the generating function must be derived from the definition of the combinatorial structure that it encodes. Then, one studies the generating function as an complex analytic function to derive asymptotic formulas of its coefficients.

Compared to analytic combinatorics in one variable, to which the transfer theorems A and A' belong, analytic combinatorics in several variables is a much less mature theory. The difference is especially stark when it comes to derving coefficient asymptotics from the generating functions (i.e. the second step outlined above). This process is often known as *singularity analysis*, since the asymptotic expansion of the coefficients of a generating function is mostly determined by the properties of the function near its singularities. For instance, if a function A(z) has a unique dominant singularity at $\rho \in \mathbb{C}^*$, then thanks to the analyticity of A everywhere else inside and on the circle of radius $|\rho|$, we can deform the contour of integration in the Cauchy integral formula (27) to a curve C that coincides with a circle of radius $r > |\rho|$ everywhere except in a neighborhood \mathcal{N} of ρ . By splitting the parts of the contour inside and outside \mathcal{N} , we get $[z^n]A(z) = I_{loc} + I_{rem}$ with

$$I_{\text{loc}} = \frac{1}{2\pi i} \int_{\mathcal{C} \cap \mathcal{N}} \frac{A(z)}{z^{n+1}} dz \quad \text{and} \quad I_{\text{rem}} = \frac{1}{2\pi i} \int_{\mathcal{C} \setminus \mathcal{N}} \frac{A(z)}{z^{n+1}} dz \,. \tag{95}$$

Since |z| = r on $\mathcal{C} \setminus \mathcal{N}$, we have $I_{\text{rem}} = O(r^{-n})$, which is exponentially small compared to $|\rho|^{-n}$. On the other hand, $\limsup_{n\to\infty} ([z^n]A(z))^{1/n} = |\rho|^{-1}$ by the root test of radius of convergence. Hence the asymptotics of $[z^n]A(z)$ is dominated by the term I_{loc} . Since I_{loc} only depends on A(z) in an arbitrarily small neighborhood of ρ , its asymptotics can be further studied via asymptotic expansions of the function A(z) near its singularity $z = \rho$. (This is basically the beginning of a proof of Theorem A'.)

At first glance, it is not clear how the above approach of singularity analysis could be generalized to multivariate functions. Indeed, the singularities of a multivariate complex function A(z) always form a continuous set with no isolated points (c.f. Hartog's extension theorem). The crucial remark here is that in general, not all singularities of A(z) contribute to the dominant asymptotics of its coefficients. In nice cases, one can even expect to find a finite number of *contributing singularities*, so that the asymptotics of $[z^n]A(z)$ is dominanted by the values of A(z) in arbitrarily small neighborhoods of these points. In practice, one would like deform the torus T_r in the Cauchy integral formula (27) to a cycle (i.e. *d*-chain without boundary) C homologous to T_r in the domain of analyticity of the integrant $\frac{A(z)}{z^{n+1}}$, such that the denominator $|z^{n+1}|$ attains its mimimum only at a finite number of points $\rho^{(1)}, \ldots, \rho^{(m)}$ on C. By Stokes' theorem (see e.g. [9, Appendix A.2]), such a deformation does not change the value of the integral. One can then take arbitrarily small neighborhoods as in the univariate setting. This gives the decomposition $[z^n]A(z) = I_{loc}^{(1)} + \cdots + I_{loc}^{(m)} + I_{rem}$, where

$$I_{\text{loc}}^{(s)} = \left(\frac{1}{2\pi i}\right)^d \int_{\mathcal{C}^{(s)}} \frac{A(\boldsymbol{z})}{\boldsymbol{z}^{\boldsymbol{n}+1}} d\boldsymbol{z} \quad \text{and} \quad I_{\text{rem}} = \left(\frac{1}{2\pi i}\right)^d \int_{\mathcal{C} \setminus \left(\mathcal{C}^{(1)} \cup \dots \cup \mathcal{C}^{(m)}\right)} \frac{A(\boldsymbol{z})}{\boldsymbol{z}^{\boldsymbol{n}+1}} d\boldsymbol{z}, \quad (96)$$

with $\mathcal{C}^{(s)} := \mathcal{C} \cap \mathcal{N}^{(s)}$. For a well chosen cycle \mathcal{C} , we expect the *non-local* term I_{rem} to be exponentially small compared to the *local* terms $I_{\text{loc}}^{(s)}$, so that the latter dominates the the asymptotics of $[\boldsymbol{z}^n]A(\boldsymbol{z})$. After this, one still needs to expand the localized integrals $I_{\text{loc}}^{(s)}$ to obtain a simple asymptotic expansion of $[\boldsymbol{z}^n]A(\boldsymbol{z})$. How this can be done depends on the form of the function $A(\boldsymbol{z})$. But usually this is a simpler problem, since there are a lot more tools at our disposal for the local analysis of a function.

The above general strategy to the multivariate singularity analysis has been outlined by Pemantle and Wilson in [9]. Over the past twenty years, they and their collaborators developped an impressive theory that treats the singularity analysis of multivariate rational-type functions in an algorithmic way. The results of their project are collected at http://acsvproject.com/. More precisely, they consider

general rational functions in several variables (and also some other functions whose singularities form algebraic varieties) and compute their coefficient asymptotics in the diagonal limit (i.e. θ -diagonal limit with $\theta = (1, ..., 1)$). In this broad setting, they carry out the strategy described in the previous paragraph in a systematic way, identifying the contributing singularities $\rho^{(s)}$ and simplifying the localized integrals $I_{loc}^{(s)}$, using powerful tools from algebraic topology and Morse theory. The reason that they need such advanced tools is that the singular set of a general rational function can have a quite complicated geometry. In particular, the location of the dominant singularities $\rho^{(1)}, \ldots, \rho^{(m)}$ (or *contributing critical points*, as they are called in [9]) depends on the direction λ of the diagonal limit. And in general, they cannot be reached by the easy-to-visualize cycles of product form $\mathcal{C} = \mathcal{C}_1 \times \cdots \times \mathcal{C}_d$.

This paper explores the case of Δ -analytic generating functions, following the same general strategy of singularity analysis. As shown by Proposition 5, this case is essentially disjoint from that of rational functions. Δ -analytic functions have a simpler singularity structure, in the sense that they always have a unique, fixed dominant singularity reachable by a cycle of product form. This allows us to study their coefficients in the θ -diagonal limits for general $\theta \in \mathbb{R}^d_{>0}$, while using mostly elementary tools from univariate complex analysis. We will discuss more the relation between our case and the case of rational functions in Section 4. For more background on ACSV, we refer to the historical accounts in [9, Chapter 1] and [8, Section 1.2.2].

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