Analytic Bezout Identities*

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1. INTRODUCTION

In a number of papers [2-4, 6, 7, 12] the problem of finding explicit solutions h_1, \ldots, h_m for the Bezout equation: $f_1h_1 + \cdots + f_mh_m = 1$ has been considered. If f_1, \ldots, f_m are complex polynomials in *n* variables and they have no common zeros in \mathbb{C}^n , the existence of explicit analytic expressions for the corresponding polynomials h_1, \ldots, h_m has a number of applications to systems theory and commutative algebra.

For instance, the problem of finding a closed loop controller for certain distributed parameter systems reduces to the question of finding a matrix with polynomial entries which is a left inverse to a rectangular matrix of polynomials of maximal rank (see [10, 20, or 21] for details). It is easy to see that this problem reduces to solving the polynomial Bezout equation (see [3]). The case where the rank is not maximal has also considerable interest in systems theory, we refer to [4] for some open questions in this case.

In several contexts, for instance, in transcendental number theory, one is interested in finding the solutions h_1, \ldots, h_m of the algebraic Bezout equation with the smallest possible degrees. Up to recently the best estimate known for deg h_j was (of the order of magnitude of) (max deg $f_j)^{2^n}$. Using explicit analytic expressions for the h_j , Brownawell [11] has shown that one can find solutions h_i with deg $h_i \leq n^2$ (max deg $f_j)^n$. This estimate is known

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to be close to optimal, as we explain in Section 3 below. Up to date no purely algebraic proof of this bound for the degrees of the h_j has been found. (The interested reader will find a short survey of this topic in [5].)

Similarly, in the deconvolution problem, one has functions f_1, \ldots, f_m which are the Fourier transforms of a strongly coprime family μ_1, \ldots, μ_m of distributions of compact support. One searches for a procedure to compute explicitly distributions of compact support ν_1, \ldots, ν_m such that $\mu_1 * \nu_1$ $+\cdots + \mu_m * \nu_m = \delta$. (Here $\hat{\nu}_j$ play the role of h_j in the equivalent formulation $\hat{\mu}_1 \hat{\nu}_1 + \cdots + \hat{\mu}_m \hat{\nu}_m \equiv 1$). This question arises in problems of robust filtering, image processing, etc. [10]. In [7] we wrote down formulas for a solution v_1, \ldots, v_m of the deconvolution problem in terms of interpolation series. The problem we have faced until recently is that, while for the one-dimensional case these formulas can be easily implemented, in the higher dimensional case they are far too cumbersome. Some of them seem to be beyond the range of symbolic languages like MACSYMA upon which we had, perhaps too optimistically, relied. For that reason we present here a new version of our original deconvolution formulas which assumes extra conditions on the family μ_1, \ldots, μ_m but has as a payoff a very simple formula for the deconvolutors v_1, \ldots, v_m . We give herein simple examples where these extra conditions are satisfied.

The problem of finding an efficient algorithm to compute the above-mentioned solutions to the algebraic Bezout equation being still open, we also analyze here the particular case in which those polynomials can be computed in terms of interpolation formulas. Finding an algorithm with a low complexity for this problem will have many important applications in the theory of distributed parameter systems and in robotics.

We have also found that a language barrier prevented our work [7] from being more easily available to some engineers, and we hope that the present paper will overcome those shortcomings.

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2. Analytic Case

We will consider only entire functions f of n complex variables satisfying inequalities of the form

$$|f(z)| \le A(1+|z|)^m d^{H(\operatorname{Im} z)}, \qquad z \in \mathbb{C}^n, \tag{1}$$

Im $z = (\text{Im } z_1, \dots, \text{Im } z_n) \in \mathbb{R}^n$, where *H* is a convex continuous function in \mathbb{R}^n , homogeneous of degree 1 (i.e., $H(\lambda x) = \lambda H(x)$ when $\lambda > 0$). We call such a function *H* a supporting function. By the Paley-Wiener theorem [18] there is a distribution μ of compact support in \mathbb{R}^n such that $f = \hat{\mu}$, the Fourier transform of μ . Furthermore, the supporting function H_0 of cv supp μ will satisfy $H_0 \leq H$ (Here cv denotes the convex hull). Conversely, if $f = \hat{\mu}$ we can take $H = H_0$ in (1) and *m* is related to the order of μ in a simple manner. Hereafter we will just write $f \in \hat{\mathscr{E}'} (= \hat{\mathscr{E}'}(\mathbb{R}^n))$ if *f* satisfies (1).

For simplicity denote $\rho(z) := \log(2 + |z|) + |\text{Im } z|$. A family f_1, \ldots, f_m of functions in $\hat{\mathscr{E}}'$ is said to be *strongly coprime* if there is a constant c such that

$$\sum_{j=1}^{m} \left| f_m(z) \right|^2 \ge e^{-c\rho(z)}, \qquad z \in \mathbb{C}^n.$$
(2)

It is well known [14] that (2) is a necessary and sufficient condition for the existence of functions $h_1, \ldots, h_m \in \hat{\mathscr{E}}'$ such that

$$\sum_{j=1}^{m} f_j h_j = 1.$$
 (3)

In other words, a strongly coprime family is precisely a family for which the analytic Bezout equation (3) has a solution. If we consider (3) in terms of the distributions $\mu_1, \ldots, \mu_m, \nu_1, \ldots, \nu_m$ such that $\hat{\mu}_j = f_j$ and $\hat{\nu}_j = h_j$ then we have the identity

$$\mu_1 * \nu_1 + \cdots + \mu_m * \nu_m = \delta, \qquad (4)$$

i.e., ν_1, \ldots, ν_m solve the deconvolution problem stated in [7]. We will say sometimes that the family of distributions μ_1, \ldots, μ_m is strongly coprime.

It might be useful to explain why (4) is called a deconvolution problem. If we have an unknown signal (function or even distribution or random process) φ then the usual data one measures would be ψ_1, \ldots, ψ_m given by

$$\psi_1 \coloneqq \mu_1 * \varphi, \dots, \psi_m \coloneqq \mu_m * \varphi. \tag{5}$$

The way to recover φ is by *deconvolution* (which is still given here by convolution with distributions or compact support).

$$\varphi = \nu_1 * \psi_1 + \cdots + \nu_m * \psi_m. \tag{6}$$

As we have mentioned in the Introduction our problem is to find easily computable functions h_j and corresponding distributions ν_j solving (3) and (4), respectively. We note that under the strongly coprime condition (2), or

even under the weaker assumption that (2) is only satisfied for real values of $z \ (z \in \mathbb{R}^n)$, there are readily available tempered distributions α_j solving the deconvolution problem, namely, let

$$\hat{\alpha}_j(z) = \frac{\bar{f}_j(z)}{\sum\limits_{j=1}^m \left|f_j(z)\right|^2}, \quad z \in \mathbb{R}^n.$$
(7)

The problem is that the α_j do not have compact support and furthermore, the α_j themselves are not so readily computable (except by inverting the Fourier transform). Nevertheless there are many situations where these α_j are still very useful, among other reasons because they minimize the noise amplification of the deconvolution process (6) (see [7] for an example of implementation in two dimensions). On the other hand, in many applications it is often not necessary to obtain an exact solution to (4) but one is allowed to replace the Dirac δ in (4) by a (sufficiently) smooth function uwith small support, i.e., an approximation to δ . It is this approximate deconvolution problem that is more readily solvable, even with very good knowledge on the support of the distributions v_j , which will turn out to be (reasonably) smooth functions.

Since the method we use relies on Koppelman-type formulas, like those developed in [1, 9], we need the following explicit relation whose proof is an immediate verification.

LEMMA 1. Let μ be a distribution of compact support in \mathbb{R}^n , $1 \le k \le n$ and $\zeta = (\zeta_1, \ldots, \zeta_n) \in \mathbb{C}^n$. The holomorphic function of 2n complex variables $g_k(z, \zeta)$ defined by

$$g_{k}(z,\zeta) := \frac{\hat{\mu}(z_{1},\ldots,z_{k},\zeta_{k+1},\ldots,\zeta_{n}) - \hat{\mu}(z_{1},\ldots,z_{k-1},\zeta_{k},\ldots,\zeta_{n})}{z_{k}-\zeta_{k}}$$
(8)

is the Fourier transform (for ζ fixed) of the distribution denoted $I = I(\mu, \zeta, k)$, which evaluated at $\varphi \in C_0^{\infty}(\mathbb{R}^n)$ has the value

$$\langle I, \varphi \rangle \coloneqq -i \int \left[\int_0^{t_k} \varphi(t_1, \dots, t_{k-1}, u, 0, \dots, 0) e^{i\zeta_k(u-t_k)} du \right]$$

$$\times e^{-i(t_{k+1}\zeta_{k+1} + \dots + t_n\zeta_n)} d\mu(t).$$
(9)

(By abuse of language we have written $\int \psi(t) d\mu(t)$ to denote $\langle \mu, \psi \rangle$).

We note that for the distribution I we have $\operatorname{cv} \operatorname{supp} \mu \subseteq \operatorname{cv} \operatorname{supp} \mu$. Furthermore the collection of functions g_1, \ldots, g_n satisfies

$$g_1(z,\zeta)(z_1-\zeta_1) + \cdots + g_n(z,\zeta)(z_n-\zeta_n) = \hat{\mu}(z) - \hat{\mu}(\zeta). \quad (10)$$

Associated to these holomorphic functions we have a (1, 0) differential form g in the variable ζ given by

$$g = g(z,\zeta,\mu) := \sum_{k=1}^{n} g_k(z,\zeta) d\zeta_k.$$
(11)

Given a family of *m* entire holomorphic functions f_1, \ldots, f_m its zero set \mathscr{Z} is defined as

$$\mathscr{Z} := \left\{ z \in \mathbb{C}^n : f_1(z) = \cdots = f_m(z) = 0 \right\}.$$
(12)

In our applications we will only consider the case where the set \mathscr{Z} is discrete. We say that \mathscr{Z} is *almost real* if there is constant A > 0 such that

$$\mathscr{Z} \subseteq \left\{ z \in \mathbb{C}^n : |\operatorname{Im} z| \le A \log(2 + |z|) \right\}.$$

It is well known that an almost real zero set \mathscr{Z} is discrete [8, 15]. For a \mathscr{Z} discrete set \mathscr{Z} , r > 0, we can define a *counting function* $n(\mathscr{Z}, r) :=$ $\#(\mathscr{Z} \cap \mathscr{Z}B_r)$, $B_r = \{z \in \mathbb{C}^n : |z| < r\} =$ Euclidean ball of center 0 and raidus r. The distance function is $d(z, \mathscr{Z}) := \min(1, \min\{|z - \zeta|: \zeta \in \mathscr{Z}\})$.

Given a family of *n* distributions of compact support in \mathbb{R}^n , μ_1, \ldots, μ_n let us denote H_1 the supporting function of $\operatorname{cv} \bigcup_{i=1}^n \operatorname{supp} \mu_i$, that is,

$$H_1(\theta) := \max_{1 \le j \le n} \max\{x \cdot \theta \colon x \in \operatorname{supp} \mu_j\}, \quad (\theta \in \mathbb{R}^n), \quad (13)$$

 $x \cdot \theta = x_1 \theta_1 + \cdots + x_n \theta_n.$

DEFINITION 1. A family of *n* distributions μ_1, \ldots, μ_n of compact support in \mathbb{R}^n is well behaved if there exist positive constants *A*, *B*, *N*, κ , a supporting function H_0 such that $0 \le H_0 \le H_1$, such that the zero set \mathscr{Z} of the functions $f_1 = \hat{\mu}_1, \ldots, f_n = \hat{\mu}_n$, is almost real,

$$n(\mathscr{Z}, r) = O(r^{A}), \tag{14}$$

and, denoting

$$|f(z)| \coloneqq \left[\sum_{1}^{n} |f_{j}(z)|^{2}\right]^{1/2},$$
(15)

the following inequality holds:

$$|f(z)| \geq \frac{Bd(z,\mathscr{Z})^{\kappa} e^{H_0(\operatorname{Im} z)}}{(1+|z|)^N}.$$
(16)

DEFINITION 2. A well-behaved family μ_1, \ldots, μ_n is very well behaved if there are constants $c_1, M, c_1 > 0$, such that for every $\zeta \in \mathscr{Z}$, we have

$$|J(\zeta)| \coloneqq \left| \det \left[\frac{\partial f_j}{\partial z_i}(\zeta) \right]_{i,j} \right| \ge c_1 (1 + |\zeta|)^{-M}.$$
(17)

This condition implies that the common zeros of f_1, \ldots, f_n are simple, that we can take $\kappa = 1$ in (16), and that if $\zeta, \zeta' \in \mathscr{Z}, \zeta \neq \zeta'$ then $|\zeta - \zeta'| \ge c_2(1 + |\zeta|)^{-M'}$ for some positive constants c_2, M' .

We will say also that functions f_1, \ldots, f_n are (very) well behaved if the above properties hold.

Given a family f_1, \ldots, f_m in $\mathscr{E}'(\mathbb{R}^n)$, m > n, with no common zeros, we introduce the following functions and differential forms. First, let $g^j = g^j(z, \zeta, \mu_j)$, $f_j = \hat{\mu}_j$, be the (1,0) differential forms in ζ given by (11), we write $g^j = \sum_{k=1}^n g_k^j d\zeta_k$. Recall the coefficients g_k^j are holomorphic in both z and ζ . Let F be the vector-valued holomorphic function $F := (f_1, \ldots, f_m)$; we write

$$|F(\zeta)| := \left[\sum_{j=1}^{m} \left|f_j(\zeta)\right|^2\right]^{1/2},$$

which is a nowhere vanishing C^{∞} function of ζ . Let

$$\varphi = \varphi(z,\zeta) \coloneqq \left[\sum_{j=1}^{m} \tilde{f}_j(\zeta) f_j(z) \right] / F(\zeta)^2, \qquad (18)$$

$$Q = Q(z,\zeta) \coloneqq \left[\sum_{j=1}^{m} \bar{f}_j(\zeta) g^j(z,\zeta) \right] / F(\zeta)^2.$$
(19)

Therefore φ is a C^{∞} function of (z, ζ) , $\varphi(\zeta, \zeta) \equiv 1$ and, as a function of z, φ is a linear combination of the f_j . Q is a (1,0) differential form in ζ and its coefficients are C^{∞} in (z, ζ) and holomorphic in z. Finally, the n + 1 functions Δ_j, C^{∞} in (z, ζ) and holomorphic in z are defined by the identities

$$g^1 \wedge \cdots \wedge g^j \wedge \cdots \wedge g^n \wedge Q = \Delta_j d\zeta_1 \wedge \cdots \wedge d\zeta_n, \quad 1 \le j \le n$$
 (20)

$$g^1 \wedge \cdots \wedge g^n = \Delta_{n+1} d\zeta_1 \wedge \cdots \wedge d\zeta_n. \tag{21}$$

It is clear that the Δ_j are simply $n \times n$ determinants whole entries are obtained from the coefficients of g^1, \ldots, g^n, Q . Therefore, as functions of z, they are finite linear combinations of products of n among the functions g'_k , $1 \le k \le n$. Note that these products are just Fourier transforms of convolutions of n distributions of the form $I(\mu_j, \zeta, k)$ (see (9)).

In order to obtain simple and easily computable deconvolution formulas we need to assume that a strongly coprime family of distributions μ_1, \ldots, μ_m contains a (very) well-behaved subfamily μ_1, \ldots, μ_n . Furthermore, we need some control on the relation between the support of all the μ_j versus the supports of the first *n*. Let

$$H_2(\theta) = \max_{1 \le j \le n} \max\{x \cdot \theta \colon x \in \operatorname{supp} \mu_j\} \qquad (\theta \in \mathbb{R}^n).$$
(22)

One such relation between the supporting functions H_0 , H_1 , H_2 is given by

$$H_2 \le 2H_1, \tag{23}$$

and

$$2(n-1)H_1(\theta) + H_2(\theta) < 2nH_0(\theta) \quad \text{if} \quad \theta \neq 0.$$
 (24)

The last condition is equivalent to

$$\exists r_0 > 0 \quad \text{such that} \quad r_0|\theta| \le 2nH_0(\theta) - 2(n-1)H_1(\theta) - H_2(\theta).$$
(25)

With all this notation in place we are now ready to state the first deconvolution formula.

THEOREM 1. Let μ_1, \ldots, μ_m be a strongly coprime family of distributions such that μ_1, \ldots, μ_n is a very well-behaved subfamily. Assume further that (23) and (25) hold. For any $u \in C_0^{\infty}(\mathbb{R}^n)$ with supp $u \subseteq \{x \in \mathbb{R}^n : |x| \le r_0\}$ one can write

$$\hat{u}(z) = \sum_{\zeta \in \mathscr{X}} \hat{u}(\zeta) \frac{\Delta_{n+1}(z,\zeta)}{J(\zeta)} \varphi(z,\zeta) + \sum_{j=1}^{n} (-1)^{N+1-j} f_j(z) \sum_{\zeta \in \mathscr{X}} \frac{\Delta_j(z,\zeta)}{J(\zeta)} \hat{u}(\zeta).$$
(26)

Formula (26) can be rewritten as

$$\hat{u}(z) = \sum_{j=1}^{m} h_j(z) f_j(z),$$

where the h_j are given by explicit interpolation formulas and they are Fourier transforms of a series of distributions which are computable in terms of the original μ_1, \ldots, μ_m . In the particular case where m = n + 1 then formula (26) can be also rewritten as

$$\hat{u}(z) = \sum_{\zeta \in Z} \frac{\hat{u}(\zeta)}{J(\zeta) f_{n+1}(\zeta)} \begin{vmatrix} g_1^1(z,\zeta) \cdots g_1^{n+1}(z,\zeta) \\ \vdots & \vdots \\ g_n^1(z,\zeta) \cdots g_n^{n+1}(z,\zeta) \\ f_1(z) \cdots f_{n+1}(z) \end{vmatrix}.$$
(27)

Proof of Theorem 1. It follows the lines of Theorem 3 from our paper [7]. It uses the Koppelman-type generalization of the Cauchy integral representation formula, especially in the version due to Anderson and Berndtsson [1]. First one introduces a parameter $\epsilon > 0$, a function φ_{ϵ} , and two (1,0) differential forms in ζ as

$$\varphi_{\varepsilon}(z,\zeta) \coloneqq \frac{\sum_{j=1}^{n} \bar{f}_{j}(\zeta) f_{j}(z) + \varepsilon}{\left| f(\zeta) \right|^{2} + \varepsilon}$$
(28)

$$s(z,\zeta) \coloneqq \sum_{j=1}^{n} \left(\bar{\zeta}_{j} - \bar{z}_{j} \right) d\zeta_{j}$$
⁽²⁹⁾

$$Q_{\epsilon}(z,\zeta) \coloneqq \frac{\sum_{j=1}^{n} \bar{f}_{j}(\zeta) g_{j}(z,\zeta)}{\left|f(\zeta)\right|^{2} + \varepsilon},$$
(30)

where, as before, $f = (f_1, \ldots, f_n) |f(\zeta)|^2 = \sum_{j=1}^n |f_j(\zeta)|^2$. The procedure from [9, pp. 402, 409]) gives two kernels K_e , P_e (i.e., differential forms in the variable ζ of type (n, n-1) and (n, n), respectively) such that if v is a holomorphic function in a neighborhood of B_R , $z \in B_R$, then

$$v(z) = \frac{1}{(2\pi i)^n} \left\{ \int_{\partial B_R} v(\zeta) K_{\epsilon}(z,\zeta) + \int_{B_R} v(\zeta) P_{\epsilon}(z,\zeta) \right\}.$$
(31)

These two kernels are defined as follows. Let $G_1(t) = t^n$ and $G_2(t) = t$; we denote for any $\alpha \in \mathbb{N}$.

$$G_1^{(\alpha)} = G_1^{(\alpha)}(z,\zeta) \coloneqq \frac{d^{\alpha}}{dt^{\alpha}} G_1|_{t=\varphi_{\epsilon}(z,\zeta)}$$
(32)

$$G_{2}^{(\alpha)} = G_{2}^{(\alpha)}(z,\zeta) := \frac{d^{\alpha}}{dt^{\alpha}} G_{2}|_{t=\varphi_{*}(z,\zeta)},$$
(33)

where φ is given by (18) and φ_e by (28). With Q_e defined by (30) and Q by

(19), we define

$$K_{\epsilon}(z,\zeta) = \sum_{\alpha_{0}+\alpha_{1}+\alpha_{2}=n-1} \frac{1}{\alpha_{1}!\alpha_{2}!} \frac{G_{1}^{(\alpha_{1})}G_{2}^{(\alpha_{2})}s \wedge (\overline{\partial}s)^{\alpha_{0}} \wedge (\overline{\partial}Q_{\epsilon})^{\alpha_{1}} \wedge (\overline{\partial}Q)^{\alpha_{2}}}{|z-\zeta|^{2(\alpha_{0}+\alpha_{1})}},$$
(34)

$$P_{e}(z,\zeta) := \sum_{\alpha_{1}+\alpha_{2}=n} \frac{1}{\alpha_{1}!\alpha_{2}!} G_{1}^{(\alpha_{1})} G_{2}^{(\alpha_{2})} (\bar{\partial}Q_{e})^{\alpha_{1}} \wedge (\bar{\partial}Q)^{\alpha_{2}}, \qquad (35)$$

where $\alpha_0, \alpha_1, \alpha_2 \in \mathbb{N}$. Everywhere the variable z is considered as a parameter and the $\overline{\partial}$ derivative is taken with respect to ζ . Due to our choice of function G_2 , the index α_2 can only take the values 0 and 1. For this reason the expression for P_e becomes particularly simple

$$P_{e} = \varphi \left(\bar{\partial} Q_{e} \right)^{n} + n \varphi_{e} \left(\bar{\partial} Q_{e} \right)^{n-1} \wedge \bar{\partial} Q.$$
(36)

The terms $(\bar{\partial}Q_{\ell})^{n-1}$ and $(\bar{\partial}Q_{\ell})^n$ must be computed, for instance,

$$\left(\overline{\partial}Q_{\epsilon}\right)^{n} = \overline{\partial}Q_{\epsilon} \wedge \cdots \wedge \overline{\partial}Q_{\epsilon} \text{ (n times)}$$
$$= (2i)^{n} n! \Delta_{n+1} \frac{\overline{J}(\zeta)\epsilon}{\left(\left|f(\zeta)\right|^{2} + \epsilon\right)^{n+1}} d\lambda, \qquad (37)$$

where $d\lambda = d\lambda(\zeta)$ = Lebesgue measure in \mathbb{C}^n . (We have eliminated the variables (z, ζ) where they were evident; we will use this convention freely in the rest of the paper.)

It is clear that φ_{ϵ} and Q_{ϵ} are singular when $\epsilon = 0$ precisely at the points $\zeta \in \mathscr{D}$. The expression (37) shows that the strength of this singularity is in one of the terms of P_{ϵ} . The strategy of the proof is to try to get very singular terms so that when $\epsilon \to 0$ the volume integrals in (31) become sums, while the boundary integrals tend to zero when we set $v = \hat{u}$ and let $R \to \infty$ over a conveniently chosen sequence. The reason this idea works is the following lemma [7, Corollary 4.1.1].

LEMMA 2. Let σ_0 be the measure which is the sum of direct masses at the points of \mathscr{X} ; i.e., for $\psi \in C_0^{\infty}(\mathbb{C}^n)$ we have $\int \psi d\sigma_0 = \sum_{\zeta \in \mathscr{X}} \psi(\zeta)$. Then, the family of measures σ_e given by

$$d\sigma_{\varepsilon}(\zeta) = \frac{\varepsilon}{\left(\left|f(\zeta)\right|^{2} + \varepsilon\right)^{n+1}} d\lambda(\zeta)$$
(38)

converges when $\varepsilon \to 0$, to the measure

$$\frac{\pi^n}{n!} \frac{d\sigma_0}{|J(\zeta)|^2},\tag{39}$$

where, as always, J denotes the determinant Jacobian of f_1, \ldots, f_n .

From (37) we see that the first term in (36) is amenable to Lemma 2. The second term is not singular enough, therefore it will be transformed using Stokes' formula in the corresponding integral of (31). Namely, due to type considerations, one obtains the first part of the identity

$$d_{\zeta} \Big\{ v(\zeta) \varphi_{\varepsilon}(z,\zeta) \Big(\left(\bar{\partial} Q_{\varepsilon}(z,\zeta) \right)^{n-1} \wedge Q(z\zeta) \Big\} \\ = \bar{\partial} \Big\{ v \varphi_{\varepsilon} \big(\bar{\partial} Q_{\varepsilon} \big)^{n-1} \wedge Q \Big\} \\ = v \varphi_{\varepsilon} \big(\bar{\partial} Q_{\varepsilon} \big)^{n-1} \wedge \bar{\partial} Q + v \bar{\partial} \varphi_{\varepsilon} \wedge \big(\bar{\partial} Q_{\varepsilon} \big)^{n-1} \wedge Q \Big\}$$

The last identity follows from the fact that v is a holomorphic function in ζ and the (2n-2) form $(\overline{\partial}Q_{\epsilon})^{n-1}$ is $\overline{\partial}$ closed. Using this identity the representation formula (31) becomes

$$v(z) = \frac{1}{(2\pi i)^{n}} \int_{\partial B_{R}} v(\zeta) \Big\{ K_{\epsilon} + n\varphi_{\epsilon} (\bar{\partial}Q_{\epsilon})^{n-1} \wedge Q \Big\} \\ + \frac{1}{(2\pi i)^{n}} \int_{B_{R}} v(\zeta) \Big\{ \varphi(\bar{\partial}Q_{\epsilon})^{n} - n\bar{\partial}\varphi_{\epsilon} \wedge (\bar{\partial}Q_{\epsilon})^{n-1} \wedge Q \Big\}, \quad (40)$$

where the integration is in the variable ζ and we have suppressed the dependency on (z, ζ) of the kernels.

LEMMA 3. The following identity holds

$$\overline{\partial}\varphi_{\epsilon} \wedge \left(\overline{\partial}Q_{\epsilon}\right)^{n-1} = (n-1)^{(n+1)n/2} \frac{\epsilon}{\left(\left|f(\zeta)\right|^{2} + \epsilon\right)^{n+1}}$$

$$\times \left[\sum_{j=1}^{n} (-1)^{j} \left(f_{j}(z) - f_{j}(\zeta)\right) \bigwedge_{k=1}^{n} g^{k}\right] \wedge \bigwedge_{k=1}^{n} \overline{\partial}f_{j}(\zeta),$$
(41)

where the wedge products in (41) are to be taken in their natural order, e.g., $\bigwedge_{k \neq 1} g^k = g^2 \wedge \cdots \wedge g^n$.

Proof of Lemma 3. We start by rewriting φ_{e} ,

$$\varphi_{\varepsilon} = \frac{\sum_{j=1}^{n} f_{j}(z) \bar{f}_{j}(\zeta) + \varepsilon}{\left|f(\zeta)\right|^{2} + \varepsilon} = \frac{\sum_{j=1}^{n} \bar{f}_{j}(\zeta) \left(f_{j}(z) - f_{j}(\zeta)\right) + \left|f_{j}(\zeta)\right|^{2} + \varepsilon}{\left|f(\zeta)\right|^{2} + \varepsilon}$$
$$= 1 + \sum_{j=1}^{n} \left(f_{j}(z) - f_{j}(\zeta)\right) \frac{\bar{f}_{j}(\zeta)}{\left|f(\zeta)\right|^{2} + \varepsilon}.$$

Denote $f_j = f_j(\zeta)$ and $\psi_j = \psi_j(\zeta) := (|f(\zeta)|^2 + \varepsilon)^{-1} \tilde{f_j}$. Then we have

$$\varphi_{\epsilon} = 1 + \sum_{j=1}^{n} (f_j(z) - f_j) \psi_j, \qquad Q_{\epsilon} = \sum_{j=1}^{n} \psi_j g^j.$$

Therefore,

$$\overline{\partial}\varphi_{\varepsilon} \wedge \left(\overline{\partial}Q_{\varepsilon}\right)^{n-1} = \left[\sum_{j=1}^{n} \left(f_{j}(z) - f_{j}\right)\overline{\partial}\varphi_{j}\right] \wedge \left[\sum_{k\neq j}\overline{\partial}\varphi_{j} \wedge g^{j}\right]^{n-1} \\ = \sum_{j=1}^{n} \left(f_{j}(z) - f_{j}\right) \left[\overline{\partial}\psi_{j} \wedge \left[\sum_{k\neq j}\overline{\partial}\psi_{k} \wedge g^{k}\right]^{n-1}\right],$$

since $\bar{\partial}\psi_j \wedge \bar{\partial}\psi_j = 0$. Using that the 2-forms $\bar{\partial}\psi_k \wedge g^k$ commute and that the product of two of them with the same index vanishes, we have

$$\left[\sum_{k\neq j}\overline{\partial}\psi_k \wedge g^k\right]^{n-1} = (n-1)! \bigwedge_{k\neq j} \left(\overline{\partial}\psi_k \wedge g^k\right)$$
$$= (n-1)! (-1)^{(n-1)n/2} \left(\bigwedge_{k\neq j} g^k\right) \wedge \left(\bigwedge_{k\neq j}\overline{\partial}\psi_k\right).$$

Hence

$$\overline{\partial}\psi_j\wedge\left[\sum_{k\neq j}\overline{\partial}\psi_k\wedge g^k\right]^{n-1}=(n-1)!(-1)^{n(n+1)/2}\left(\bigwedge_{k\neq j}g^k\right)\wedge\left(\bigwedge_{k=1}\overline{\partial}\psi_k\right).$$

Now, we have $\bar{\partial}\psi_k = (|f|^2 + \varepsilon)^{-1} \overline{\partial} f_k - (|f|^2 + \varepsilon)^{-2} \bar{f}_k \overline{\partial} |f|^2$. Therefore we can use that $\bar{\partial} |f|^2 \wedge \bar{\partial} |f|^2 = 0$ and obtain

$$\bigwedge_{k=1}^{n} \overline{\partial} \psi_{k} = \frac{\bigwedge_{k=1}^{n} \overline{\partial} f_{k}}{\left(|f|^{2} + \varepsilon\right)^{n}} - \frac{\sum_{j=1}^{n} \left[f_{j} \bigwedge_{1 \le n < j}^{n} \overline{\partial} f_{k} \wedge \overline{\partial} |f|^{2} \wedge \bigwedge_{j < k \le n}^{n} \overline{\partial} f_{k} \right]}{\left(|f|^{2} + \varepsilon\right)^{n+1}}$$

If we now expand $\overline{\partial}|f|^2 = \sum_k f_k \overline{\partial} f_k$, we see that only the term $f_j \overline{\partial} f_j$ remains in the triple product above. Hence

$$\bigwedge_{k=1}^{n} \overline{\partial} \psi_{k} = \frac{\bigwedge_{k=1}^{n} \overline{\partial} f_{k}}{\left(|f|^{2} + \varepsilon\right)^{n}} - \frac{\left\lfloor \sum_{j=1}^{n} |f_{j}|^{2} \right\rfloor \bigwedge_{k=1}^{n} \overline{\partial} f_{k}}{\left(|f|^{2} + \varepsilon\right)^{n+1}} = \frac{\varepsilon}{\left(|f|^{2} + \varepsilon\right)^{n+1}} \bigwedge_{k=1}^{n} \overline{\partial} f_{k}.$$

This concludes the proof of Lemma 3. \Box

Lemma 3 tells us that $\overline{\partial}\varphi_{\epsilon} \wedge (\overline{\partial}Q_{\epsilon})^{n-1} \wedge Q$ is the product of a measure with a smooth density, independent of ϵ , and the function $\epsilon(|f(\zeta)|^2 + \epsilon)^{n-1}$. Lemma 2 can now be invoked to see that the volume integral in (40) reduces to a sum when $\epsilon \to 0$. In fact, let us choose R so that $|f(\zeta)|^2 =$ $|f_1(\zeta)|^2 + \cdots + |f_n(\zeta)|^2 \neq 0$ when $|\zeta| = R$. This choice is always possible since \mathscr{X} is a discrete set by assumption. In this case none of φ_{ϵ} , K_{ϵ} , Q_{ϵ} have singularities when $\epsilon = 0$ and $\zeta \in \partial B_R$. We set φ_0 , K_0 , Q_0 to be the correspondent quantities. Therefore

$$v(z) = \frac{1}{(2\pi i)^{n}} \int_{\partial B_{R}} v(\zeta) \Big\{ K_{0} + n\varphi_{0} \Big((\bar{\partial}Q_{0})^{n-1} \wedge Q \Big) \Big\} \\ + \frac{1}{(2\pi i)^{n}} \lim_{\epsilon \to 0^{+}} \int_{B_{R}} v(\zeta) \Big[\varphi(\bar{\partial}Q_{\epsilon})^{n} - n\bar{\partial}\varphi_{\epsilon} \wedge (\bar{\partial}Q_{\epsilon})^{n-1} \wedge Q \Big].$$

Recall that $f_j(\zeta) = 0$ if $\zeta \in \mathscr{Z}$ and $1 \le j \le n$. Using Lemma 2 and the definition (20) and (21) of the Δ_j we can compute explicitly the limit and obtain

$$v(z) = \frac{1}{(2\pi i)^n} \int_{\partial B_R} v(\zeta) \Big(K_0 + n\varphi_0 \Big((\overline{\partial} Q_0)^{n-1} \wedge Q \Big) \Big) \\ + \sum_{\zeta \in \mathscr{Z} \cap B_R} v(\zeta) \frac{\Delta_{n+1}(z,\zeta)\varphi(z,\zeta)}{J(\zeta)} + \sum_{j=1}^n (-1)^{n+1-j} f_j(z) \\ \times \sum_{\zeta \in \mathscr{Z} \cap B_R} v(\zeta) \frac{\Delta_j(z,\zeta)}{J(\zeta)}.$$
(42)

Up to this moment we have only used that $\mathscr{Z} \cap \partial B_R = \emptyset$ and that v is holomorphic in a neighborhood of \overline{B}_R . To let $R \to \infty$ we have to choose a sequence $R \to \infty$ judiciously. Recall that $n(\mathscr{Z}, r) = \#(\mathscr{Z} \cap B_r) \leq Cr^A$ for some positive constants A, D, and all $r \geq 1$. Let M be the smallest integer $\geq C(R+1)^A + 1$, divide the shell $\overline{B}_{R+1} \setminus B_R$ into M concentric subshells but choosing the boundaries to be $\partial B_{(R+j/M)}$, $0 \le j \le M$. There is at least one such subshell that is free from points of \mathscr{D} ; choose R' to be the mid-radius of this subshell, then $d(\zeta, \mathscr{Z}) \ge (2M)^{-1}$ if $|\zeta| = R'$. Starting from the sequence R = q = 1, 2, ... we construct a sequence $R_q, q < R_q <$ q + 1, such that for some positive constants A_1, N_1 ,

$$|f(\zeta)| \ge A_1 q^{-N_1} e^{H_0(\operatorname{Im} \zeta)} \quad \text{if} \quad |\zeta| = R_q.$$
 (43)

This follows from (16) and the choice of R_q . We are now ready to estimate the terms in the boundary integral of (42) for $R = R_f$. We will assume $|z| \le C_0 < \infty$ and consider those q such that $R_q \ge C_0 + 1.$

First, let us observe that the functions g_k^j , $1 \le j \le m$, $1 \le k \le n$, satisfy an estimate of the form

$$\left|g_{k}^{j}(z,\zeta)\right| \leq C_{1}(1+|z|)^{M_{1}}(1+|\zeta|)^{M_{1}}e^{H_{2}(\operatorname{Im} z)+H_{2}(\operatorname{Im} \zeta)}, \quad (44)$$

for some constants $M_1, C_1 > 0$. If $1 \le j \le n$, we can replace H_2 by H_1 . We can now estimate the coefficients of differential form Q. Denote $||Q(z, \zeta)||$ the largest absolute value of the coefficients of $d\zeta_k$ at the point (z, ζ) . We proceed as follows. First,

$$|F(\zeta)| \ge |f(\zeta)| \ge A_1 q^{-M_1} e^{H_0(\operatorname{Im} \zeta)} \quad \text{if } |\zeta| = R_q.$$
(45)

Therefore,

$$\begin{aligned} \|Q(z,\zeta)\| &\leq \frac{1}{|F(\zeta)|} \left[\sum_{j=1}^{n} \|g^{j}(z,\zeta)\|^{2} \right]^{1/2} \\ &\leq \frac{C_{1}(1+|z|)^{M_{1}} e^{H_{2}(\operatorname{Im} z)} (1+R_{q})^{M_{1}} e^{H_{2}(\operatorname{Im} z)}}{A_{1}q^{-N_{1}} e^{H_{0}(\operatorname{Im} \zeta)}} \end{aligned}$$

which leads to

$$\|Q(z,\zeta)\| \le C_3 q^{N_2} e^{H_2(\operatorname{Im} \zeta) - H_0(\operatorname{Im} \zeta)}.$$
(46)

The constant C_3 depends in fact on z, but $|z| \le C_0$ and $N_2 = M_1 + N_1$. (In fact C_3 can be estimated in terms of $e^{2H_2(\text{Im }z)}$ and polynomials in |z|.) Similarly, with possibly different values for the constants C_3 , N_2 appearing below, we have

$$\left\|\bar{\partial}Q(z,\zeta)\right\| \le C_3 q^{N_2} e^{H_2(\operatorname{Im} z) - 2H_0(\operatorname{Im} z)} \tag{47}$$

$$\|\bar{\partial}Q_0(z,\zeta)\| \le C_3 q^{N_2} e^{2H_1(\operatorname{Im} z) - 2H_0(\operatorname{Im} z)}$$
(48)

$$|\varphi_0(z,\zeta)| \le C_3 q^{N_1} e^{-H_0(\operatorname{Im} z)}$$
 (49)

$$|\varphi(z,\zeta)| \leq \frac{|F(z)|}{|F(\zeta)|} \leq C_3 q^{N_1} e^{-H_0(\operatorname{Im} z)}.$$
 (50)

To estimate K_0 we recall that α_2 can only take the values 0 and 1 in (34). In case $\alpha_2 = 0$ we have to estimate terms of the form $\varphi_0^{n-\alpha_1}\varphi(\bar{\partial}Q_0)^{\alpha_1}$, with $0 \le \alpha_1 \le n-1$. There are powers of q that we will disregard, the estimate is then the function of Im ζ ,

$$\left\|\varphi_0^{n-\alpha}\varphi(\bar{\partial}Q_0)^{\alpha_1}\right\| \ll e^{-(n+\alpha_1+1)H_0+2\alpha_1H_1}$$

Since $H_1 \ge H_0$, the worst case estimate occurs when $\alpha_1 = n - 1$. Hence the terms corresponding to $\alpha_2 = 0$, $\alpha_0 + \alpha_1 = n - 1$, in the definition of K_0 can all be estimated by

$$C_{3}q^{N_{3}}e^{-2nH_{0}(\operatorname{Im}\zeta)+2(n-1)H_{1}(\operatorname{Im}z)}, \qquad |\zeta| = R_{q}.$$
 (51)

The terms with $\alpha_1 = 1$, $\alpha_0 + \alpha_1 = n - 2$ corresponds to the estimate of $\|\varphi_0^{n-\alpha_1}(\bar{\partial}Q_0)^{\alpha_1} \wedge \bar{\partial}Q\|$. The worst case occurs this time when $\alpha_1 = n - 2$ and we obtain an estimate of the form

$$C_{3}q^{N_{3}}e^{-2nH_{0}(\operatorname{Im}\zeta)+2(n-2)H_{1}(\operatorname{Im}\zeta)+2H_{2}(\operatorname{Im}\zeta)}, \qquad |\zeta| = R_{q}.$$
(52)

In (42) we have one more term to estimate for $|\zeta| = R_a$,

$$\left\|\varphi_{0}\left(\bar{\partial}Q_{0}\right)^{n-1}\wedge Q\right\| \leq C_{3}q^{N_{3}}e^{-2nH_{0}(\operatorname{Im}\zeta)+2(n-1)H_{1}(\operatorname{Im}\zeta)+H_{2}(\operatorname{Im}\zeta)}.$$
 (53)

The conditions (23) and (25) imply that the largest exponential factor in (51), (52), and (53) is the one in (53) and it satisfies

$$-2nH_0(\operatorname{Im}\zeta) + 2(n-1)H_1(\operatorname{Im}\zeta) + H_2(\operatorname{Im}\zeta) \le -r_0|\operatorname{Im}\zeta|.$$
(54)

Since we have assumed that $u \in C_0^{\infty}(\overline{B}_{r_0})$, we have

$$\hat{u}(\zeta) | \le C_4 (1 + |\zeta|)^{-N_3 - 2n} e^{r_0 |\operatorname{Im}\zeta|}, \quad \zeta \in \mathbb{C}^n,$$
 (55)

which allows us to conclude that, with $v = \hat{u}$,

$$\lim_{q\to\infty}\int_{\partial B_{R_q}}\hat{u}(\zeta)\Big[K_0(z,\zeta)+n\varphi_0(z,\zeta)\big(\bar{\partial}Q_0\big)^{n-1}\wedge Q(z,\zeta)\Big]=0.$$
 (56)

To conclude the proof of Theorem 1 we need to show that the series appearing in the representation formula (26) converge absolutely and uniformly in compact subsets of \mathbb{C}^n to functions in \mathscr{E}' . Since we have assumed that \mathscr{Z} is almost real, the estimates of all the terms Δ_j , φ , J^{-1} are in terms of powers of $|\zeta|$. Using that $n(\mathscr{Z}, r) = O(r^A)$ and $|\hat{u}(\zeta)|$ decreases as fast as $|\zeta|^{-N_4}$ for $\zeta \in \mathscr{Z}$, we have the desired convergence once N_4 is chosen sufficiently large. The support of all the distributions thus obtained is contained in the convex set K whose supporting function is $(n + 1)H_2$. \Box

Remarks 1. One can see that the condition $2H_1 \ge H_2$ cannot be relaxed if the other conditions of the theorem remain the same, otherwise the exponent in (52) would become positive and we would not be able to prove (56).

2. A way to weaken the conditions in Theorem 1 is to impose some better lower bound on |F| than (45) that only depends on the first *n* functions. We will do so in Theorem 2 below.

3. It is clear that one only needs $u \in C_0^N(\overline{B}_{r_0})$ for N sufficiently large to obtain (26).

The following example shows how Theorem 1 enormously simplifies the computation of the deconvolution formula proposed in [7].

Let μ_1, μ_2, μ_3 be the characteristic functions of the squares centered at 0, of sides parallel to the axes and of lengths $2\sqrt{3}, 2\sqrt{2}, 2$, respectively. One can easily show [7, 17] that μ_1, μ_2 is a very well-behaved family with

$$H_0(\operatorname{Im}\zeta) = \sqrt{2} \left(|\operatorname{Im}\zeta_1| + |\operatorname{Im}\zeta_2| \right)$$

Here $H_1(\operatorname{Im} \zeta) = H_2(\operatorname{Im} \zeta) = \sqrt{3} (|\operatorname{Im} \zeta_1| + |\operatorname{Im} \zeta_2|)$. In this case the main hypotheses (24) reduces to verify that $4\sqrt{2} - 3\sqrt{3} > 0!$. Since $|x_1| + |x_2| = \sqrt{x_1^2 + x_2^2}$ for $x \in \mathbb{R}^2$, we have $r_0 = 4\sqrt{2} - 3\sqrt{3} \ge 0.2$. The variety \mathscr{Z} in this case is given by

$$\mathscr{Z} = \left\{ \left[\frac{j\pi}{\sqrt{3}}, \frac{k\pi}{\sqrt{2}} \right] : j, k \in \mathbb{Z}^* \right\} \cup \left\{ \left[\frac{j\pi}{\sqrt{2}}, \frac{k\pi}{\sqrt{3}} \right] : j, k \in \mathbb{Z}^* \right\}.$$

(There were about forty different types of terms in [7] to compute.)

Before we proceed to state Theorem 2 we need to point out that the representation formula (31) does not depend on the particular choice of the differential forms g^{j} we have chosen, rather on the fact that (10) is satisfied. That is,

$$\sum_{k=1}^{n} g_{k}^{j}(z,\zeta)(z_{k}-\zeta_{k}) = f_{j}(z) - f_{j}(\zeta).$$

Now, let $f = \hat{\mu}$ and $h(z) = \sin Bz_1$ for some B > 0, and denote by g the differential forms associated to f by (8) and (11). Let us define a differential form γ by

$$\gamma(z,\zeta) := f(\zeta) \frac{\sin Bz_1 - \sin B\zeta_1}{z_1 - \zeta_1} d\zeta_1 + h(z)g(z,\zeta), \qquad (57)$$

writing $\gamma = \sum \gamma_k d\zeta_k$ we have

$$\gamma_1(z,\zeta)(z_1-\zeta_1)+\cdots+\gamma_n(z,\zeta)(z_n-\zeta_n)=f(z)h(z)-f(\zeta)h(\zeta).$$

Therefore we can associate γ to the product $f \cdot h$. It is also clear that as a function of z, the γ_k are Fourier transforms of distribution of compact support, easily computable terms of μ and B. Obviously we can replace sin Bz_1 by sin Bz_j without any problems; hence, given a family f_1, \ldots, f_m we can construct an *augmented family* $f_1, \ldots, f_m, f_{m+1} \coloneqq f_1 \sin Bz_1, \ldots, f_{2m}$ $\coloneqq f_m \cdot \sin Bz_1, \ldots, f_{(n+1)m} \coloneqq f_m \cdot \sin Bz_n$. The corresponding g^j for $j \ge m + 1$ are computed following the procedure (57). It is clear that if f_1, \ldots, f_m was strongly coprime, the augmented family we will keep the notation H_0, H_1, H_2 to indicate the support functions corresponding to the m original members of the augmented family $f_1, \ldots, f_{(n+1)m}$.

THEOREM 2. Let f_1, \ldots, f_m be a strongly coprime family such that the subfamily f_1, \ldots, f_n is very well behaved. There are constants $B_0 \ge 0, r_0 > 0$, such that for any $B \ge B_0$, and any $u \in C_0^{\infty}(\overline{B}_{r_0})$, the representation formula (26) is valid for the augmented family $f_1, \ldots, f_m, \ldots, f_{(n+1)m}$ defined above if either of the following two conditions holds:

$$H_2 \le 2H_1$$
 and $2(n-1)H_1 < (2n-1)H_0$ (58)

$$2H_1 < H_2$$
 and $2(n-2)H_1 + H_2 < (2n-1)H_0$. (59)

Proof. The proof is exactly the same as that of Theorem 1 except for improvements on the estimates (46), (47), and (50) for the new Q and φ . Recall that it is there where all the functions $f_1, \ldots, f_{(n+1)m}$ appear. Let $F_1 := (f_1, \ldots, f_{(n+1)m})$ and keep the notation $F = (f_1, \ldots, f_m)$ as before. We have

$$|F_1(\zeta)| = |F(\zeta)| (1 + |\sin B\zeta_1|^2 + \dots + |\sin B\zeta_n|^2)^{1/2} = |F(\zeta)| \cdot \vartheta(\zeta).$$
(60)

It is clear that for some positive constants c_n, c'_n we have

$$\vartheta(\zeta) = \left(1 + |\sin B\zeta_1|^2 + \dots + |\sin B\zeta_n|^2\right)^{1/2}$$

$$\geq c'_n \left[e^{2B |\operatorname{Im} \zeta_1|} + \dots + e^{2B |\operatorname{Im} \zeta_n|}\right]^{1/2}$$

$$\geq c_n e^{2(B/n) |\operatorname{Im} \zeta|_1} \geq c_n^{2(B/n) |\operatorname{Im} \zeta|},$$

where $|\operatorname{Im} \zeta|_1 = \sum_{j=1}^n |\operatorname{Im} \zeta_j|$, $|\operatorname{Im} \zeta| = [\sum_{j=1}^n |\operatorname{Im} \zeta_j|^2]^{1/2}$. We estimate first g^j , j = m + 1, for $|z - \zeta| \ge 1$, since that is the only case that appears in the proof of (56). As it follows from (57) we have, for some *i*, k $(1 \le i \le n, 1 \le k \le m)$, the estimate

$$\begin{split} \|g^{j}(z,\zeta)\| &\leq c e^{2B |\operatorname{Im} z_{i}|} \Big[\|f_{k}(\zeta)| (1+|\sin B\zeta_{i}|^{2})^{1/2} + \|g^{k}(z,\zeta)\| \Big] \\ &\leq C(1+|z|)^{N_{1}} (1+|\zeta|)^{N_{1}} e^{B |\operatorname{Im} z_{i}|} \\ &\times \Big[\|f(\zeta)\|\vartheta(\zeta) + e^{H_{2}(\operatorname{Im} z)} e^{H_{2}(\operatorname{Im} \zeta)} \Big]. \end{split}$$

It follows that

$$\begin{split} \|Q(z,\zeta)\| &\leq \frac{1}{|F_1(\zeta)|} \left[\sum_{j=1}^{(n+1)m} \|g^j(z,\zeta)\|^2 \right]^{1/2} \\ &\leq C_1 (1+|z|)^{N_1} (1+|\zeta|)^{N_1} e^{B|\operatorname{Im} z|+H_2(\operatorname{Im} z)} \\ &\times \left\{ \frac{e^{H_2(\operatorname{Im} \zeta)}}{|F_1(\zeta)|} + \frac{|F(\zeta)|\vartheta(\zeta)}{|F_1(\zeta)|} \right\}. \end{split}$$

Therefore, for a positive constant C_2 depending on z, we obtain

$$\|Q(z,\zeta)\| \le C_2 (1+|\zeta|)^{N_1} [1+e^{H_2(\operatorname{Im}\zeta)-H_0(\operatorname{Im}\zeta)-(B/n)|\operatorname{Im}\zeta|_1}].$$
(61)

Similarly,

$$|\varphi(z,\zeta)| \leq \frac{|F_1(z)|}{|F_1(\zeta)|} \leq C_2(1+|\zeta|)^{N_1} e^{H_0(\operatorname{Im}\zeta) - (B/n)|\operatorname{Im}\zeta|_1}.$$

Finally,

$$\|\overline{\partial}Q(z,\zeta)\| \leq C \max_{1\leq j, \ k\leq (n+1)m} \frac{\|\partial f_j(\zeta)\| \|g^k(z,\zeta)\|}{|F_1(\zeta)|^2}.$$

Now, for $j \ge m + 1$ we have for some $1 \le i \le n, 1 \le l \le m$,

$$\begin{aligned} \left\| \partial f_j(\zeta) \right\| &= \left\| \sin B\zeta_1 \cdot \partial f_i(\zeta) + f_i(\zeta) B \cdot \cos B\zeta_i \, d\zeta_i \right\| \\ &\leq C_1 (1 + |\zeta|)^{N_2} \big(\vartheta(\zeta) e^{H_2(\operatorname{Im} \zeta)} + |F(\zeta)| \vartheta(\zeta) \big). \end{aligned}$$

It follows that

Choose $B_0 \ge 0$ such that

$$2H_2(\operatorname{Im} \zeta) - 2H_0(\operatorname{Im} \zeta) + (B_0/n)|\operatorname{Im} \zeta|_1 \le 0.$$
 (63)

When $B \ge B_0$ we will have $||Q(z,\zeta)|| \le C_2(1+|\zeta|)^{N_1}$ and

$$\left\|\bar{\partial}Q(z,\zeta)\right\| \le C_2(1+|\zeta|)^{N_1} e^{H_2(\operatorname{Im}\zeta)-H_0(\operatorname{Im}\zeta)},\tag{64}$$

where C_2 still denotes a constant depending on z of the form

$$C_2 = \text{const}(1 + |z|)^{N_1} e^{H_2(\text{Im } z) + B |\text{Im } z|_1}.$$

We can now return to the proof of Theorem 1 at the point where we obtained the estimates (51), (52), (53). Ignoring powers of q, the exponential factors are

$$\exp(-2nH_0(\mathrm{Im}\,\zeta) + 2(n-1)H_1(\mathrm{Im}\,\zeta) - (B/n)|\mathrm{Im}\,\zeta|_1) \quad (51')$$

$$\exp(-(2n-1)H_0(\operatorname{Im}\zeta) + 2(n-2)H_1(\operatorname{Im}\zeta) + H_2(\operatorname{Im}\zeta)) \quad (52')$$

$$\exp(-(2n-1)H_0(\operatorname{Im}\zeta) + 2(n-2)H_1(\operatorname{Im}\zeta)). \tag{53'}$$

Under hypothesis (58) the largests of these three is (53') and its exponent satisfies

$$-(2n-1)H_0(\operatorname{Im} \zeta) + 2(n-1)H_1(\operatorname{Im} \zeta) \le -r_0|\operatorname{Im} \zeta|, \quad (53'')$$

for some $r_0 > 0$. If the hypothesis (59) holds, then the largest exponent is

(52') and we define $r_0 > 0$ by

$$-(2n-1)H_0(\operatorname{Im}\zeta) + 2(n-2)H_1(\operatorname{Im}\zeta) + H_2(\operatorname{Im}\zeta) \le -r_0|\operatorname{Im}\zeta|.$$
(52")

In either case the rest of the proof is the same as that of Theorem 1.

EXAMPLE. As [7] shows the family μ_1, μ_2, μ_3 obtained by taking $\mu_1 =$ characteristic function of the unit square = $\chi_{[-1,1]\times[-1,1]}, \mu_2$ a rotation of μ_1 by 36° and μ_3 a rotation of μ_1 by 45° satisfies the first conditions of Theorem 1 and Theorem 2 with $H_0(\theta) = |\theta|$, since the squares contain the unit disk. One can easily convince oneself that the hypothesis (24) does not hold (e.g., take $\theta = (t, 0), t > 0$.) On the other hand, one can take $H_1(\theta) = H_2(\theta) = \sqrt{2} |\theta|$, since all the squares are contained in the disk of radius $\sqrt{2}$. We are in the situation of hypothesis (58) and its verification reduces to the fact that

$$r_0 = 3 - 2\sqrt{2} > 0.$$

Furthermore, $B_0 = 4(\sqrt{2} - 1)$ works in this case.

3. POLYNOMIAL CASE

The conditions on Theorems 1 and 2 imply that the convex set defined by H_0 contains a ball. If we want to prove an algebraic version of (26), the fact that this condition is not satisfied plays a role. Such a representation was stated in [3, 4] without proof. We analyze here the conditions under which it is valid.

THEOREM 3. Let p_1, \ldots, p_m be a family of polynomials in \mathbb{C}^n without common zeros, suppose further that:

- (a) $D := \max_{1 \le i \le m} \deg p_i = \deg p_i$ for $1 \le i \le n$.
- (b) $\mathscr{L} = \{ z \in \mathbb{C}^n : p_i(z) = 0; 1 \le i \le n \}$ is discrete.

(c) J(z) := Jacobian determinant of p_1, \ldots, p_n at z is $\neq 0$ for all $z \in \mathscr{Z}$.

(d) p_1, \ldots, p_n have no common zeros at infinity, i.e., $\#\mathscr{L} = D^n$.

Then

$$1 = \sum_{\boldsymbol{\zeta} \in \mathscr{Z}} \frac{\Delta_{n+1}(\boldsymbol{z},\boldsymbol{\zeta})}{J(\boldsymbol{\zeta})} \varphi(\boldsymbol{z},\boldsymbol{\zeta}) + \sum_{j=1}^{n} (-1)^{n+1-j} p_j(\boldsymbol{z}) \sum_{\boldsymbol{\zeta} \in \mathscr{Z}} \frac{\Delta_j(\boldsymbol{z},\boldsymbol{\zeta})}{J(\boldsymbol{\zeta})}, \quad (65)$$

where φ, Δ_j are defined as in (18)–(21) with respect to the polynomials p_1, \ldots, p_m .

Remarks 1. The functions g_k defined by (8) are obviously polynomials of degree D - 1. It follows that (66) has the form

$$p_1(z)A_1(z) + \dots + p_m(z)A_m(z) = 1$$
 (66)

for some polynomials $A_j \in \mathbb{C}[z_1, \ldots, z_n]$ of degree of most n(D-1). This follows from the fact they are given as $n \times n$ determinants involving the g_k^j .

2. As before, the case m = n + 1 leads to a particularly pleasing form of (66),

$$1 = \sum \frac{1}{J(\zeta)} \frac{1}{p_{n+1}(\zeta)} \begin{vmatrix} g_1^1(z,\zeta) \cdots g_1^{n+1}(z,\zeta) \\ \vdots & \vdots \\ g_n^1(z,\zeta) \cdots g_n^{n+1}(z,\zeta) \\ p_1(z) \cdots p_{n+1}(z) \end{vmatrix}.$$
(67)

3. The two statements in condition (d) above are really a form of Bezout's theorem [19]. The meaning of the expression " p_1, \ldots, p_n have no common zeros at infinity" is that if we introduce homogeneous polynomials $H_j(z_0, \ldots, z_n) := z_0^D p_j[z_1/z_0, \ldots, z_n/z_0]$ then the subset of \mathbb{C}^{n+1} defined by $\{z_0 = 0, H_1 = \cdots = H_n = 0\}$ is $\{0\}$. This is equivalent to the statement

$$|p(z)| = \left[\sum_{j=1}^{n} |p_j(z)|^2\right]^{1/2} \ge C|z|^D \quad \text{if } |z| \ge R_0 \ge 1.$$
 (68)

If we call $p_1^0(z)$ the leading homogeneous polynomial of p_j , then the estimate (68) is also equivalent to

$$\left\{z \in \mathbb{C}^n : p_1^0(z) = \cdots = p_n^0(z) = 0\right\} = \{0\}.$$
 (69)

We also note that (69) implies (b) above. That is, condition (d) above implies condition (b).

Proof of Theorem 3. The proof is the same as that of Theorem 1. This time we take v = 1 and R arbitrary $\geq R_0$ (cf. (68)). One can estimate $\|Q(z, \zeta)\| \leq C/|\zeta|$, $\|\overline{\partial}Q(z, \zeta)\| \leq C/|\zeta|^2$, $|\varphi_0(z, \zeta)| \leq C/|\zeta|^2$, $\|\varphi_0(z, \zeta)\| \leq C/|\zeta|^2$, $\|\varphi_0(z, \zeta)\| \leq C/|\zeta|^2$ if $|\zeta| = R$ and $|z| \leq K \leq R - 1$, with C = C(k) > 0. These estimates imply that the boundary integral in (40) tends to zero when $R \to \infty$. \Box

If p_1, \ldots, p_n are such that their leading terms p_j^0 satisfy (69) but their degrees deg $p_j = D_j$ are not all equal or max{ $D_j: 1 \le j \le n$ } is smaller than $D = \max\{D_j: 1 \le j \le m\}$, then we can still prove a version of Theorem 3. That corresponds to the analytic counterpart of Theorem 1, that is, to

Theorem 2. For the moment we continue to assume that $J(z) \neq 0 \ \forall z \in \mathscr{Z}$ = { $p_1 = \cdots = p_n = 0$ }.

Let $L_1(z) = u_1 z_1 + \cdots + u_n z_n$ be a linear homogeneous polynomial with generic coefficients. The condition that $\{z \in \mathbb{C}^n : L_1 = p_2^0 = \cdots = p_2^0 = 0\}$ $\neq \{0\}$ is an algebraic condition on the coefficients of L_1 . Therefore we can choose L_1 such that for any integer $d_1 \ge 0$ we have $\{[p_1^0 L_1^{d_1} = p_2^0 = \cdots = p_2^0 = 0\} = \{0\}$. Continuing in this fashion we can choose L_1, \ldots, L_n , d_1, \ldots, d_n , such that for any choice of constants $\varepsilon_1, \ldots, \varepsilon_n$, if we define $\tilde{p}_j = p_j (L_j + \varepsilon_j)^{d_j}$, then $\{\tilde{p}_1^0 = \cdots = \tilde{p}_n^0 = 0\} = \{0\}$ and deg $\tilde{p}_j = D$ for $1 \le j \le n$. It is clear now that for most choices of ε_j we still have that all common zeros of $\tilde{p}_1^0 = \cdots = \tilde{p}_n^0$ are simple and $\tilde{p}_1^0 = \cdots = \tilde{p}_n^0, p_{n+1}, \ldots, p_m$ have no common zeros. Theorem 3 can now be applied to this new family; one obtains polynomials A_j ,

$$\sum_{j=1}^{m} A_{j} p_{j} = 1, \qquad \deg A_{j} \le n(D-1),$$
(70)

and such that they have a representation of the type (65).

We remark that a representation such as (65) cannot be valid if the p_1, \ldots, p_n have common zeros at infinity. For instance, in the example of Masser and Philippon in [11],

$$p_1 = z_1^D, p_2 = z_1 - z_2^D, \dots, p_{n-1} = z_{n-2} - z_{n-1}^D, p_n = 1 - z_{n-1} z_n^{D-1}$$

one knows that $\delta = D^n - D^{n-1}$ is the best estimate possible for the degrees of A_i solving the polynomial Bezout equation. The polynomials

$$A_{1} = z_{n}^{\delta}, A_{2} = -z_{n}^{\delta} \frac{\left[z_{1}^{D} - z_{2}^{D^{2}}\right]}{z_{1} - z_{2}^{D}}, \qquad A_{3} = -z_{n}^{\delta} \frac{\left[z_{2}^{D} - z_{3}^{D^{3}}\right]}{z_{2} - z_{3}^{D}}, \dots,$$
$$A_{n-1} = -z_{n}^{\delta} \frac{\left[z_{n-2}^{D^{n-2}} - z_{n-1}^{D^{n-1}}\right]}{z_{n-2} - z_{n-1}^{D}}, \qquad A_{n} = \frac{1 - z_{n-1}^{D^{n-1}} z_{n}^{\delta}}{1 - z_{n-1} z_{n}^{D^{-1}}},$$

have exactly this degree. On the other hand, if we had a representation like (65) we could conclude that there are solutions A_j of the polynomial Bezout equation with deg $A_j \le n(D-1)$ as in (70).

We would like now to show that the condition (c) for the simplicity of the zeros in Theorem 3 is not necessary. Regretfully, we only know how to do this in the case where m = n + 1.

THEOREM 4. Let p_1, \ldots, p_{n+1} be a fimily of polynomials in \mathbb{C}^n without any common zeros, $D = \deg p_1 = \cdots = \deg p_n \ge \deg p_{n+1}$ and $\{p_1^0 = \cdots = p_n^0 = 0\} = \{0\}$. Then we can find polynomials A_j of degree \le

n(D-1) satisfying the identity $\sum_{1}^{n+1}A_jp_j = 1$. The coefficients of the A_j can be written in terms of the values of p_{n+1} and values of the derivatives of p_{n+1} , and the coefficients of $g_k^{n+1}(z,\zeta)$ (when considered as polynomials in z), all of these evaluated at the points of $\mathscr{Z} = \{\zeta : p_1(\zeta) = \cdots = p_n(\zeta) = 0\}$.

Proof. For $0 < \epsilon_j \ll 1$ the function $1/p_{n+1}$ is holomorphic in a neighborhood of $\Pi = \{|p_1| \le \epsilon_1, \ldots, |p_n| \le \epsilon_n\}$. This set is a compact polynomially convex set. (The compactness follows from condition (68)). By Sard's theorem one can choose the ϵ_j so that the sets $\{|p_j| = \epsilon_j\}$ are real analytic submanifolds of \mathbb{C}^n . (In fact, we only need it in a neighborhood of Π .) For an $v \in \mathscr{H}(\overline{\Pi})$ we have that integral

$$\operatorname{Res}_{\mathscr{Z}}(vd\zeta_{1}\wedge\cdots\wedge d\zeta_{n}):=\frac{1}{(2\pi i)^{n}}\int_{|p_{1}|=\sigma_{1}}v(\zeta)\frac{d\zeta_{1}\wedge\cdots\wedge d\zeta_{n}}{p_{1}(\zeta)\cdots p_{n}(\zeta)} \quad (71)$$

is independent of choice of $\sigma_1, \ldots, \sigma_n$ as long as $0 < \sigma_j \le \varepsilon_j$ and the $\{|p_j| = \sigma_j\}$ are smooth. Furthermore, if v is in the ideal generated by p_1, \ldots, p_n in $\mathscr{H}(\Pi)$ then this residue is zero. Therefore, it depends only on the values of v at \mathscr{Z} and a certain number of derivatives of v at \mathscr{Z} (as it follows from the Nullstellensatz as presented, e.g., in [13, 16]). In other words, the integral (71) can be considered as an operator defined by a certain linear combination of the Dirac masses δ_{ζ} and their derivatives $(\partial^{|\alpha|}/\partial \zeta^{\alpha}) \delta_{\zeta}, \zeta \in \mathscr{Z}$, applied to the holomorphic function v. This operator is very hard to compute explicitly except in very simple cases but it is perfectly defined as the common value of all the integrals (71). It is called the *residue current* of \mathscr{Z} .

Let us consider now the polynomial B(z) defined by

$$B(z) = \operatorname{Res}_{\mathscr{Z}} \frac{1}{p_{n+1}(\zeta)} \begin{vmatrix} g_1^1(z,\zeta) \cdots g_1^{n+1}(z,\zeta) \\ \vdots & \vdots \\ g_n^1(z,\zeta) \cdots g_n^{n+1}(z,\zeta) \\ p_1(z) \cdots p_{n+1}(z) \end{vmatrix} d\zeta_1 \wedge \cdots \wedge d\zeta_n.$$
(72)

This polynomial is in fact of the form $\sum_{j=1}^{n+1} A_j(z) p_j(z)$, with A_j polynomials of degree $\leq n(D-1)$ whose coefficients are computed in terms of the values of derivatives of p_{n+1} and the coefficients of g_k^j (as polynomials in z) evaluated over \mathscr{D} . The only problem is to show $B(z) \equiv 1$. We fix values σ_j , $0 < \sigma_j < \varepsilon_j$. Consider complex numbers $\alpha_1, \ldots, \alpha_n$ sufficiently small and so chosen that:

(i) All the common zeros of $p_1 - \alpha_1, \ldots, p_n - \alpha_n$ are simple and lie in $\{|p_1| < \frac{1}{2}\sigma_1, \ldots, |p_n| < \frac{1}{2}\sigma_n\}.$

(ii) $p_1 - \alpha_1, \ldots, p_n - \alpha_n, p_{n+1}$ have no common zeros.

Note that $p_1 - \alpha_1, \ldots, p_n - \alpha_n$ still do not have any common zeros at ∞ . Let us denote $\mathscr{Z}_{\alpha} = \{z \in \mathbb{C}^n : p_1 - \alpha_1 = \cdots = p_n - \alpha_n = 0\}$. Then $\operatorname{Res}_{\mathscr{T}_{a}}(vd\zeta_{1}\wedge\cdots\wedge d\zeta_{n})$ $=\frac{1}{(2\pi i)^n}\int_{|P_1|=\sigma_1}v(\zeta)\frac{d\zeta_1\wedge\cdots\wedge d\zeta_n}{(p_1(\zeta)-\alpha_1)\cdots(p_n(\zeta)-\alpha_n)}$ $|p_n| = \sigma_n$ $= \sum_{\zeta \in \mathscr{Z}_{\alpha}} v(\zeta) / J(\zeta),$ (73) where

$$J(\zeta) = \frac{\partial(p_1 - \alpha_1, \dots, p_n - \alpha_n)}{\partial(\zeta_1, \dots, \zeta_n)} = \frac{\partial(p_1, \dots, p_n)}{\partial(\zeta_1, \dots, \zeta_n)}$$

and the last identity follows from Stokes' theroem. (Replace the contour by little spheres about the distinct points of \mathscr{Z}_{a}). This last identity is the effective computation of the residue current of \mathscr{Z}_{α} , namely,

$$\operatorname{Res}_{\mathscr{Z}_{\alpha}} = \sum_{\zeta \in \mathscr{Z}_{\alpha}} J(\zeta)^{-1} \delta_{\zeta}.$$

The first identity in (73) shows that $\operatorname{Res}_{\mathscr{Z}_{\sigma}} \to \operatorname{Res}_{\mathscr{Z}}$ as $\sigma \to 0$, i.e., the residue currents at $\alpha = 0$ are continuous when acting on holomorphic (n, 0)forms. On the other hand, by Theorem 3 we have

$$\operatorname{Res}_{\mathscr{Z}_{\alpha}}\left(\frac{1}{p_{n+1}(\zeta)}\left|\begin{array}{c}g_{1}^{1}(z,\zeta)\cdots g_{1}^{n+1}(z,\zeta)\\g_{n}^{1}(z,\zeta)\cdots g_{n}^{n+1}(z,\zeta)\\p_{1}(z)-\alpha_{1}\cdots p_{n+1}(z)\end{array}\right|d\zeta_{1}\wedge\cdots\wedge d\zeta_{n}\right)\equiv 1.$$

(Note that the g^{j} corresponding to p_{j} and to $p_{j} - \alpha_{j}$ coincide.) By continuity we obtain $B \equiv 1$. This concludes the proof of Theorem 4. \Box

Remarks 1. We can obviously obtain the same result without assuming the degrees of p_1, \ldots, p_n coincide or that they are larger or equal than that of p_{n+1} .

2. The reasoning of Theorem 4 extends to a strongly coprime family of n + 1 elements whose first n members form a well-behaved family. Under the other conditions of Theorem 1 or Theorem 2, we obtain a series representation of the solutions of Bezout equation which we computed in terms of the residue current associated to \mathcal{Z} . This time the series converges after grouping of terms.

3. The interest of the theorems in this section lies in the search for explicit algorithms to obtain solutions A_i for the algebraic Bezout equation $\sum_{i=1}^{m} A_{i} p_{i} = 1$ which satisfy Brownawell's estimate, deg $A_{i} \leq 3\mu n D^{\mu}$, $\mu =$ $\min(n, m)$.

4. CONCLUSION

We have shown how explicit solutions to the analytic and algebraic Bezout equations can be obtained under natural restrictions on the original functions f_1, \ldots, f_m . This work has applications to the implementation of deconvolution for multidetector systems.

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