# CONVERGENCE OF TWO-DIMENSIONAL WEIGHTED INTEGRALS 

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#### Abstract

A two-dimensional weighted integral in $\mathbb{R}^{2}$ is proposed as a tool for analyzing higher-dimensional unweighted integrals, and a necessary and sufficient condition for the finiteness of the weighted integral is obtained.


## 1. Introduction

In this paper, we consider weighted integrals in $\mathbb{R}^{2}$ of the form

$$
\begin{equation*}
\int_{B} \frac{|g(x, y)|^{\epsilon}}{|f(x, y)|^{\delta}} d y d x \tag{1.1}
\end{equation*}
$$

where $f$ and $g$ are real-analytic (possibly complex-valued) functions on $\mathbb{R}^{2}, \epsilon, \delta$ are positive numbers and $B$ is a small ball centered at the origin. The goal of the paper is to answer the question of finiteness for the above integral.

A simpler version of the above problem has a rich history in the literature, where it arises in connection with the growth rate of real analytic functions and decay rate of oscillatory integrals. In particular, for $g \equiv 1$ and in dimension $n=2$, the problem of finiteness of the integral in (1.1) has been fully treated by Phong, Stein and Sturm [2], while the related problem of determining the oscillation index of a two-dimensional oscillatory integral dates back to Varchenko [5]. It would be of considerable interest, however, to tackle the issue in higher dimensions (i.e. $n \geqslant 3$ ), where the problem of convergence of the unweighted integral $(g \equiv 1)$ is still poorly understood and finiteness results are, at best, partial (see [5]).

Weighted integrals of the form (1.1) sometimes arise from their unweighted higher dimensional analogues after a suitable change of coordinates, especially if the higher dimensional $f$ comes equipped with certain symmetries that can be exploited to reduce the dimension. A case in point is the important counterexample to Arnold's problem given by Varchenko in the context of oscillatory integrals in $\mathbb{R}^{3}$ (see section 5 of [5]). In our situation, Varchenko's example translates to

$$
\begin{equation*}
\iiint_{B^{3} \subset \mathbb{R}^{3}} \frac{d x_{1} d x_{2} d x_{3}}{\left|\left(\lambda x_{1}^{2}+x_{1}^{4}+x_{2}^{2}+x_{3}^{2}\right)^{2}+x_{1}^{4 p}+x_{2}^{4 p}+x_{3}^{4 p}\right|^{\delta}}, \tag{1.2}
\end{equation*}
$$

where $\lambda$ is a real parameter and $p$ is a sufficiently large natural number. Now, a few trivial size estimates coupled with a cylindrical change of coordinates transforms

[^0]the above three-dimensional unweighted integral to a two-dimensional weighted one, given by
\[

$$
\begin{equation*}
\iint_{B^{2} \subset \mathbb{R}^{2}} \frac{|y| d y d x}{\left|\left(y^{2}+x^{4}+\lambda x^{2}\right)^{2}+x^{4 p}+y^{4 p}\right|^{\delta}} \tag{1.3}
\end{equation*}
$$

\]

Clearly, the integral in (1.2) converges if and only if the integral in (1.3) does. In general, the hope is that results for integrals of the form (1.1) would shed some light on the behavior of the higher-dimensional unweighted ones they arise from.

Another problem closely related to the unweighted integrals is the one of obtaining $L^{2}$ bounds for oscillatory integral operators with degenerate real-analytic phase. The case for $n=2$ has been addressed by Phong and Stein in [1]. In section 3 , the techniques used in their paper are adapted to settle the problem of local integrability of $|g|^{\epsilon} /|f|^{\delta}$.

The problem of finiteness of the integral in (1.1) also gives rise to several related problems, some of which in turn translate to results for unweighted integrals in higher dimensions. A few questions which arise naturally in this context are:
(a) When is the ratio $|g|^{\epsilon} /|f|^{\delta}$ an $A_{p}$ weight in $B$ ?
(b) How does the integral in (1.1) behave under perturbations of $f$ if $g$ is held fixed?

We shall return to these questions in subsequent papers.

## 2. Notation, Definitions and Preliminary Observations

Let $f$ and $g$ be real-analytic functions in a neighborhood of the origin in $\mathbb{R}^{2}$ such that $f(0)=g(0)=0, f, g \not \equiv 0$. Suppose the Taylor series expansions of $f$ and $g$ are given by

$$
f(x, y)=\sum_{\substack{p, q=0 \\ p, q) \neq(0,0)}}^{\infty} c_{p q} x^{p} y^{q}, \quad g(x, y)=\sum_{\substack{p, q=0 \\(p, q) \neq(0,0)}}^{\infty} \tilde{c}_{p q} x^{p} y^{q}
$$

The Newton polyhedron of $f(x, y)$ (respectively $g(x, y)$ ) at the origin is defined to be the convex hull of the union of all the north-east quadrants in $\mathbb{R}_{\geqslant 0}^{2}$ with corners at the points $(p, q)$ satisfying $c_{p q} \neq 0$ (respectively $\left.\tilde{c}_{p q} \neq 0\right)$. The boundary of the Newton polyhedron is called the Newton diagram.

There is an alternative description of the Newton diagram which is sometimes more useful for analytical purposes. In view of the Weierstrass Preparation Theorem, $f$ can be expressed, after a non-singular change of coordinates and up to a non-vanishing factor which we ignore, as a polynomial in $y$ with coefficients analytic in $x$. Factoring out this polynomial, we write $f$ as

$$
f(x, y)=x^{\tilde{\alpha}_{1}} y^{\tilde{\beta}_{1}} \prod_{\nu ; \nu \in I_{f}}\left(y-r_{\nu}(x)\right)
$$

where $\tilde{\alpha}_{1}$ and $\tilde{\beta}_{1}$ are non-negative integers and the $r_{\nu}(x)$ 's are the non-trivial zeros of the above-mentioned polynomial in $y$. Here, $\nu$ ranges over an index set $I_{f}$ that is in one-to-one correspondence with the set of roots of $f$. In order to avoid confusion with subscripts later on, it is better to think of $\nu$ not as a positive integer but as the unique element of $I_{f}$ associated to the root $r_{\nu}$. In a small neighborhood of 0 , the $r_{\nu}$ 's may be expressed as fractional power series, the so-called Puiseux series in
$x$, of the following form:

$$
r_{\nu}(x)=c_{\nu} x^{a_{\nu}}+O\left(x^{b_{\nu}}\right)
$$

where $b_{\nu}>a_{\nu}>0$ are rational numbers and $c_{\nu} \neq 0$. In fact, each $r_{\nu}(x)$ is an analytic function in $z$, where $z=x^{1 / M}$ for some positive integer $M$. A detailed discussion of the theory of existence and convergence of Puiseux series may be obtained in Siegel [4, pp. 90-98] and in Saks and Zygmund [3, pp. 268-271]. For a more concise treatment, see section 3 of [1].

We may also write, using similar arguments,

$$
g(x, y)=x^{\tilde{\alpha}_{2}} y^{\tilde{\beta}_{2}} \prod_{\mu ; \mu \in I_{g}}\left(y-s_{\mu}(x)\right)
$$

where the $s_{\mu}$ 's are the roots of $g$, and the leading exponent of $s_{\mu}$ is denoted by $a_{\mu}$. We order the combined set of distinct exponents $a_{\nu}$ and $a_{\mu}$ into one increasing list of exponents $a_{l}$,

$$
0<a_{1}<a_{2}<\cdots<a_{N}
$$

and define

$$
m_{l}:=\#\left\{\nu: r_{\nu}(x)=c_{\nu} x^{a_{l}}+\cdots, \text { for some } c_{\nu} \neq 0\right\}
$$

We call $m_{l}$ the generalized multiplicity of $f$ corresponding to the exponent $a_{l}$. The generalized multiplicity of $g$ corresponding to $a_{l}$ will be denoted by $n_{l}$. If $a_{l}$ does not occur as a leading exponent of any root of $f$ (respectively $g$ ), then we set $m_{l}=0$ (respectively $n_{l}=0$ ).

Now, it can be shown (for a proof see observation 1 in section 5 (b) of [1]) that the Newton diagram of $f$ has vertices at the points $\left(A_{l}, B_{l}\right)$, where

$$
\begin{array}{ll}
A_{l}=\tilde{\alpha}_{1}+a_{1} m_{1}+a_{2} m_{2}+\cdots+a_{l} m_{l}, & B_{l}=\tilde{\beta}_{1}+m_{l+1}+\cdots+m_{N} \\
C_{l}=\tilde{\alpha}_{2}+a_{1} n_{1}+a_{2} n_{2}+\cdots+a_{l} n_{l}, & D_{l}=\tilde{\beta}_{2}+n_{l+1}+\cdots+n_{N}
\end{array}
$$

It follows that the leading exponents $a_{l}$ of the roots of $f$ (or $g$ ) can be read off from their respective Newton diagrams, together with their generalized multiplicities $m_{l}$. In fact, the boundary segment joining the vertices $\left(A_{l-1}, B_{l-1}\right)$ and $\left(A_{l}, B_{l}\right)$ has slope $\frac{-1}{a_{l}}$ and vertical drop $m_{l}$.

Next, we introduce the following notation. If $\mathcal{L}_{l}$ is a line of slope $\frac{-1}{a_{l}}$ that occurs in the Newton diagram of either $f$ or $g$, then we denote the point of intersection of $\mathcal{L}_{l}$ with the bisectrix $p=q$ by

$$
\begin{cases}\left(\delta_{l}^{-1}, \delta_{l}^{-1}\right) & \text { if } \mathcal{L}_{l} \text { is a boundary line of the Newton diagram of } f, \\ \left(\tilde{\delta}_{l}^{-1}, \tilde{\delta}_{l}^{-1}\right) & \text { if } \mathcal{L}_{l} \text { is a boundary line of the Newton diagram of } g\end{cases}
$$

The alternative description of the Newton diagram allows us to obtain an explicit formula for $\delta_{l}^{-1}$ and $\tilde{\delta}_{l}^{-1}$ in terms of $a_{l}, A_{l}, B_{l}, C_{l}$ and $D_{l}$ :

$$
\begin{equation*}
\delta_{l}=\frac{1+a_{l}}{A_{l}+a_{l} B_{l}}, \quad \tilde{\delta}_{l}=\frac{1+a_{l}}{C_{l}+a_{l} D_{l}} \tag{2.1}
\end{equation*}
$$

The Newton distance of $f$, denoted by $\delta_{0}(f)$, is defined as follows:

$$
\delta_{0}(f):=\min _{l} \delta_{l}
$$

The Newton distance plays an important role in problems involving the growth rate of the distribution function of $f$, local integrability of $|f|^{-\delta}$ and the decay rate of
$L^{2}$ norms of oscillatory integral operators with phase $f$. In particular, it is known that $|f|^{-\delta}$ is locally integrable in a neighborhood of the origin if $\delta<\delta_{0}(f)$ (after possibly a certain analytic change of coordinates), while the exact decay rate of an oscillatory integral operator on $\mathbb{R}$ with phase $f(x, y)$ and smooth non-degenerate cut-off $\chi(x, y)$ is given by $\delta_{0}\left(f_{x y}\right)$. In this paper, we shall obtain an analogue of $\delta_{0}(f)$ in the weighted situation.

We conclude our preliminary discussion with the following observations:
(1) First, we note that the quantities $a_{l}, m_{l}, n_{l}, A_{l}, B_{l}, C_{l}, D_{l}, \delta_{l}, \tilde{\delta}_{l}$ are all coordinate dependent.
(2) Second, we observe that the above discussion allows us to extend the notion of Newton diagram to a certain class of functions strictly larger than the class of real-analytic functions. Let $f$ (and $g$ ) be of the form $\prod_{\kappa}\left(y-t_{\kappa}(x)\right)$ (where the $t_{\kappa}$ 's are convergent Puiseux series in $x$, though not necessarily the roots of an analytic function). Then the definitions of $a_{l}, m_{l}$, etc. still make sense, with the possibility that now the vertices $\left(A_{l}, B_{l}\right)$ may be fractional, instead of integer points. Let us consider the set

$$
\bigcup_{l}\left\{(x, y) ; x \geqslant A_{l}, y \geqslant B_{l}\right\}
$$

and call the boundary of the convex hull of the above union set the generalized Newton diagram of $f$. The geometric interpretation of $\left(\delta_{l}^{-1}, \delta_{l}^{-1}\right)$ and $\left(\tilde{\delta}_{l}^{-1}, \tilde{\delta}_{l}^{-1}\right)$ continues to hold, and we compute these objects using the same formulae as in (2.1).
(3) Finally, we claim that although $\delta_{l}^{-1}$ (respectively $\tilde{\delta}_{l}^{-1}$ ) has been defined only for those values of $l$ which are "represented" in the Newton diagram of $f$ (respectively $g$ ), it is possible to define these quantities for all values of $l, 1 \leqslant l \leqslant N$, using the convention that $m_{l}=0$ (respectively $n_{l}=0$ ) for any $l$ that is "omitted" in $f$ (respectively $g$ ). Geometrically, if $m_{l}=0$, one imagines a line segment of length zero and of slope $\frac{-1}{a_{l}}$ in between the segments of slope $\frac{-1}{a_{l-1}}$ and $\frac{-1}{a_{l+1}}$ in the (generalized) Newton diagram of $f$, and defines $\left(\delta_{l}^{-1}, \delta_{l}^{-1}\right)$ to be the point of intersection of this "imaginary" line segment with the bisectrix $p=q$. The coordinates of the point of intersection are still obtained from the expressions in (2.1).

The statement of our main theorem requires the following definition:
Let us consider a coordinate transformation $\varphi$ given by

$$
(x, y) \mapsto(x, y-q(x)) \quad \text { or } \quad(x, y) \mapsto(x-q(y), y)
$$

where $q$ is a convergent real-valued Puiseux series in a neighborhood of the origin. We shall see in the proof of Theorem 1 that a change of coordinates of this type transforms a real-analytic function to a function of the form

$$
\prod_{\kappa}\left(y-t_{\kappa}(x)\right)
$$

to which the notion of a generalized Newton diagram may be applied. We shall call such a transformation $\varphi$ "good" and denote the class of all good transformations by $\mathcal{C}$. For any $\varphi \in \mathcal{C}$, we define the weighted Newton distance of $f$ and $g$ associated to $\varphi$ as follows:

$$
\delta_{0}(g, f, \epsilon ; \varphi):=\min _{l}\left[\delta_{l}(\varphi)\left(1+\frac{\epsilon}{\tilde{\delta}_{l}(\varphi)}\right)\right]
$$

where the index $l$ runs through the combined set of slopes of boundary lines of the (generalized) Newton diagrams of $f$ and $g$ expressed in the coordinate system $\varphi$. $\delta_{l}(\varphi)$ and $\tilde{\delta}_{l}(\varphi)$ are the values of $\delta_{l}$ and $\tilde{\delta}_{l}$, respectively, computed in these coordinates.

## 3. Statement and Proof of Main Theorem

We are now ready to state and prove the main theorem.
Theorem 1. Let $f$ and $g$ be real-analytic (complex-valued) functions in a neighborhood of the origin in $\mathbb{R}^{2}$ with $f(0)=g(0)=0$. Then, there exists a small neighborhood $V$ of the origin such that for every $\epsilon>0$ the following holds:

$$
\int_{V} \frac{|g(x, y)|^{\epsilon}}{|f(x, y)|^{\delta}} d y d x<\infty
$$

for any $\delta<\delta_{0}^{\mathcal{C}}(g, f, \epsilon)$, and

$$
\int_{V} \frac{|g(x, y)|^{\epsilon}}{|f(x, y)|^{\delta}} d y d x=\infty
$$

for any $\delta \geq \delta_{0}^{\mathcal{C}}(g, f, \epsilon)$, where

$$
\delta_{0}^{\mathcal{C}}(g, f, \epsilon):=\min _{\varphi \in \mathcal{C}} \delta_{0}(g, f, \epsilon ; \varphi)
$$

In fact, there exists a finite subclass of $\mathcal{C}$, denoted by $\mathcal{C}_{0}$, which depends on $f$ and $g$ but is independent of $\epsilon$, such that

$$
\delta_{0}^{\mathcal{C}}(g, f, \epsilon)=\delta_{0}^{\mathcal{C}_{0}}(g, f, \epsilon):=\min _{\varphi \in \mathcal{C}_{0}} \delta_{0}(g, f, \epsilon ; \varphi)
$$

The class $\mathcal{C}_{0}$ contains, apart from the identity transformation, coordinate changes of the form $(x, y) \mapsto(x, y-q(x))$ and $(x, y) \mapsto(x-q(y), y)$ for certain Puiseux series $q$ (depending on $f$ or $g$ ) with leading exponents greater than or equal to 1 .
Remark. We shall henceforth call any $\varphi \in \mathcal{C}_{0}$ an admissible change of coordinates. Our proof, in fact, provides an algorithm for computing all the admissible coordinate changes and for selecting the optimal one. One also concludes from the theorem that any $\varphi \in \mathcal{C}_{0}$, and in particular the "minimizing" transformation, can be at worst $C^{1+\epsilon_{0}}$ for some small $\epsilon_{0}>0$ depending on $f$ and $g$.
Proof of Theorem 1. Let us for now assume that $a_{1} \geqslant 1$. This results in no loss of generality in the proof that we are about to present, and we shall mention the routine modifications for the general case at the end. Also without loss of generality, we restrict ourselves to considering the integral only on the quadrant $\{x>0, y>0\}$, and inspect the integrand on the support of the cut-off functions $\chi_{j}(x) \chi_{k}(y)$, where $\chi_{j}(x)$ and $\chi_{k}(y)$ are the indicator functions of the sets

$$
2^{-j-1}<x<2^{-j+1} \quad \text { and } \quad 2^{-k-1}<y<2^{-k+1}
$$

respectively. Here $j$ and $k$ may be taken to be large. We need to consider the following ranges of $j$ and $k$, depending on whether there is any cancellation between $y$ and $x^{a_{l}}$ :

$$
\begin{gathered}
k \leqslant a_{1} j-K_{1} ; \quad k \geqslant a_{N} j+K_{N} \\
a_{l} j+K_{l} \leqslant k \leqslant a_{l+1} j-K_{l+1} \text { for some } l \text { with } 1 \leqslant l \leqslant N-1 \\
a_{l} j-K_{l} \leqslant k \leqslant a_{l} j+K_{l} \text { for some } l \text { with } 1 \leqslant l \leqslant N
\end{gathered}
$$

Here the $K_{l}$ 's are large constants depending only on the exponents $a_{l}$ and possibly on the coefficients of the Puiseux series in the factorization of $f$ and $g$. For convenience, we shall denote the first range by $k \ll a_{1} j$, the second one by $k \gg a_{N} j$, the third category of ranges by $a_{l} j \ll k \ll a_{l+1} j$ and the fourth by $k \approx a_{l} j$. The binary relation $A \sim B$ is used in the following sense:

$$
A \sim B \quad \text { iff } \quad \frac{1}{C_{1}} B \leq A \leq C_{2} B
$$

for large constants $C_{1}$ and $C_{2}$ depending on $f$ and $g$ alone.
Case 1: $k \ll a_{1} j$. In this case,

$$
|f(x, y)| \sim 2^{-k\left(\tilde{\beta}_{1}+m_{1}+m_{2}+\cdots+m_{N}\right)}=2^{-k B_{0}}, \quad|g(x, y)| \sim 2^{-k D_{0}}
$$

Therefore,

$$
\begin{aligned}
& \sum_{(k, j) ; k \ll a_{1} j} \int \frac{|g(x, y)|^{\epsilon}}{|f(x, y)|^{\delta}} \chi_{j}(x) \chi_{k}(y) d y d x \\
& \quad \sim \sum_{k \ll a_{1} j} 2^{k\left(\delta B_{0}-\epsilon D_{0}\right)} 2^{-k} 2^{-j}=\sum_{k \ll a_{1} j} 2^{k\left(\delta B_{0}-\epsilon D_{0}-1\right)} 2^{-j} \\
& \quad \sim \sum_{k} 2^{k\left(\delta B_{0}-\epsilon D_{0}-1-\frac{1}{a_{1}}\right)}
\end{aligned}
$$

which converges iff

$$
\delta B_{0}-\epsilon D_{0}-1-\frac{1}{a_{1}}<0
$$

i.e., iff

$$
\begin{aligned}
\delta & <\frac{\epsilon D_{0}+1+\frac{1}{a_{1}}}{B_{0}}=\frac{\epsilon a_{1} D_{0}+a_{1}+1}{a_{1} B_{0}} \\
& =\frac{\epsilon \frac{a_{1} D_{0}}{a_{1}+1}+1}{\frac{a_{1} B_{0}}{a_{1}+1}}=\delta_{1}\left(1+\frac{\epsilon}{\tilde{\delta}_{1}}\right) .
\end{aligned}
$$

The case $k \gg a_{N} j$ is treated in a similar way.
Case 2: $a_{l} j \ll k \ll a_{l+1} j$. In this case,

$$
\left|y-r_{\nu}(x)\right| \sim \begin{cases}2^{-a_{\nu} j} & \text { if } a_{\nu} \leqslant a_{l} \\ 2^{-k} & \text { if } a_{\nu}>a_{l}\end{cases}
$$

We then have the following size estimates:

$$
\begin{aligned}
|f(x, y)| & \sim 2^{-j\left(\tilde{\alpha}_{1}+a_{1} m_{1}+\ldots+a_{l} m_{l}\right)} 2^{-k\left(\tilde{\beta}_{1}+m_{l+1}+\ldots+m_{N}\right)} \\
& \sim 2^{-A_{l} j} 2^{-B_{l} k} .
\end{aligned}
$$

Similarly,

$$
|g(x, y)| \sim 2^{-C_{l} j} 2^{-D_{l} k}
$$

Therefore,

$$
\begin{aligned}
& \quad \sum_{a_{l} j \ll k \ll a_{l+1} j} \int \frac{|g(x, y)|^{\epsilon}}{|f(x, y)|^{\delta}} \chi_{j}(x) \chi_{k}(y) d y d x \\
& \quad \sim \sum_{a_{l} j \ll k \ll a_{l+1} j} 2^{j\left(\delta A_{l}-\epsilon C_{l}\right)} 2^{k\left(\delta B_{l}-\epsilon D_{l}\right)} 2^{-k} 2^{-j},
\end{aligned}
$$

where the sum is over both indices $j$ and $k$. To facilitate the summation, let us introduce the new dummy variable $\tilde{k}$ in place of $k$, defined as follows:

$$
k=\tilde{k}+a_{l} j, \quad 0 \ll \tilde{k} \ll\left(a_{l+1}-a_{l}\right) j
$$

Then the above expression reduces to

$$
\begin{aligned}
& \sum_{0 \ll \tilde{k} \ll\left(a_{l+1}-a_{l}\right) j} 2^{j\left(\delta A_{l}-\epsilon C_{l}-1\right)} 2^{\left(\tilde{k}+a_{l} j\right)\left(\delta B_{l}-\epsilon D_{l}-1\right)} \\
&=\sum_{0 \ll \tilde{k} \ll\left(a_{l+1}-a_{l}\right) j} 2^{\left[\delta\left(A_{l}+a_{l} B_{l}\right)-\epsilon\left(C_{l}+a_{l} D_{l}\right)-a_{l}-1\right] j} \times 2^{\tilde{k}\left(\delta B_{l}-\epsilon D_{l}-1\right)},
\end{aligned}
$$

which converges iff
(a) $\delta\left(A_{l}+a_{l} B_{l}\right)-\epsilon\left(C_{l}+a_{l} D_{l}\right)-\left(a_{l}+1\right)<0$ when $\delta B_{l}-\epsilon D_{l}-1<0$,
(b) $\delta\left(A_{l}+a_{l} B_{l}\right)-\epsilon\left(C_{l}+a_{l} D_{l}\right)-\left(a_{l}+1\right)+\left(a_{l+1}-a_{l}\right)\left(\delta B_{l}-\epsilon D_{l}-1\right)<0$

$$
\text { when } \delta B_{l}-\epsilon D_{l}-1 \geqslant 0 \text {; }
$$

i.e. iff
(a) $\delta<\quad \frac{\epsilon\left(C_{l}+a_{l} D_{l}\right)+a_{l}+1}{A_{l}+a_{l} B_{l}}=\delta_{l}\left(1+\frac{\epsilon}{\tilde{\delta}_{l}}\right), \quad$ when $\delta<\frac{\epsilon D_{l}+1}{B_{l}}$,
(b) $\delta<\frac{\epsilon\left(C_{l}+a_{l+1} D_{l}\right)+a_{l+1}+1}{A_{l}+a_{l+1} B_{l}}=\frac{\epsilon\left(C_{l+1}+a_{l+1} D_{l+1}\right)+a_{l+1}+1}{A_{l+1}+a_{l+1} B_{l+1}}$

$$
=\delta_{l+1}\left(1+\frac{\epsilon}{\tilde{\delta}_{l+1}}\right), \text { when } \delta \geqslant \frac{\epsilon D_{l}+1}{B_{l}} .
$$

Now, a straightforward computation shows that the following conditions are equivalent:
(i) $\frac{\epsilon D_{l}+1}{B_{l}}<\frac{\epsilon\left(C_{l}+a_{l+1} D_{l}\right)+a_{l+1}+1}{A_{l}+a_{l+1} B_{l}}=\delta_{l+1}\left(1+\frac{\epsilon}{\tilde{\delta}_{l+1}}\right)$;
(ii) $\frac{\epsilon D_{l}+1}{B_{l}}<\frac{\epsilon\left(C_{l}+a_{l} D_{l}\right)+a_{l}+1}{A_{l}+a_{l} B_{l}}=\delta_{l}\left(1+\frac{\epsilon}{\tilde{\delta}_{l}}\right)$;
(iii) $\frac{\epsilon\left(C_{l}+a_{l+1} D_{l}\right)+a_{l+1}+1}{A_{l}+a_{l+1} B_{l}}<\frac{\epsilon\left(C_{l}+a_{l} D_{l}\right)+a_{l}+1}{A_{l}+a_{l} B_{l}}$;
(iv) $\frac{\epsilon\left(C_{l}+a_{l+1} D_{l}\right)+a_{l+1}+1}{A_{l}+a_{l+1} B_{l}}<\frac{\epsilon C_{l}+1}{A_{l}}$;
(v) $\frac{\epsilon\left(C_{l}+a_{l} D_{l}\right)+a_{l}+1}{A_{l}+a_{l} B_{l}}<\frac{\epsilon C_{l}+1}{A_{l}}$;
(vi) $\frac{\epsilon D_{l}+1}{B_{l}}<\frac{\epsilon C_{l}+1}{A_{l}}$.

So, there are two possibilities:
(a)

$$
\frac{\epsilon D_{l}+1}{B_{l}}<\frac{\epsilon C_{l}+1}{A_{l}},
$$

in which case the necessary and sufficient conditions in (i) - (vi) lead to the string of inequalities

$$
\frac{\epsilon D_{l}+1}{B_{l}}<\delta_{l+1}\left(1+\frac{\epsilon}{\tilde{\delta}_{l+1}}\right)<\delta_{l}\left(1+\frac{\epsilon}{\tilde{\delta}_{l}}\right)<\frac{\epsilon C_{l}+1}{A_{l}} .
$$

Therefore the optimal $\delta$ satisfies

$$
\delta<\delta_{l+1}\left(1+\frac{\epsilon}{\tilde{\delta}_{l+1}}\right)
$$

The second possibility is

$$
\begin{equation*}
\frac{\epsilon D_{l}+1}{B_{l}} \geqslant \frac{\epsilon C_{l}+1}{A_{l}} \tag{b}
\end{equation*}
$$

in which case

$$
\frac{\epsilon D_{l}+1}{B_{l}} \geqslant \delta_{l+1}\left(1+\frac{\epsilon}{\tilde{\delta}_{l+1}}\right) \geqslant \delta_{l}\left(1+\frac{\epsilon}{\tilde{\delta}_{l}}\right) \geqslant \frac{\epsilon C_{l}+1}{A_{l}}
$$

and the optimal $\delta$ satisfies

$$
\delta<\delta_{l}\left(1+\frac{\epsilon}{\tilde{\delta}_{l}}\right)
$$

In either case, we have

$$
\delta<\min _{l}\left[\delta_{l}\left(1+\frac{\epsilon}{\tilde{\delta}_{l}}\right)\right]
$$

where $\delta_{l}$ and $\tilde{\delta}_{l}$ are computed based on $f$ and $g$ in the original set of coordinates, i.e. under the identity transformation $(x, y) \mapsto(x, y)$.

Case 3: $k \approx a_{l} j$. This range of $j$ and $k$ depicts the region "near" the cluster of roots of $f$ and $g$ that has a common leading exponent $a_{l}$. For the sake of a more organized notation to be made clear shortly, let us rename $l=l_{1}$. We now fix a real root $q_{1}$ of the above cluster, such that

$$
q_{1}(x)=c_{l_{1}}^{\alpha_{1}} x^{a_{l_{1}}}+\text { higher order terms }
$$

for some index $\alpha_{1}$ with $c_{l_{1}}^{\alpha_{1}} \neq 0$ (the case of complex-valued roots will be treated shortly). Since $q_{1}$ is real-valued for all $x$ in its domain, $c_{l_{1}}^{\alpha_{1}} \in \mathbb{R}$. Here, $q_{1}$ could be the Puiseux series of a root of $f$ and/or $g$. For a real-analytic function $h$, which in our case might be either $f$ or $g$, we introduce the following notation:

$$
\begin{aligned}
\mathcal{S}_{l_{1}}(h) & :=\{q ; q \text { is the Puiseux series corresponding to a root of } h \\
& \left.\quad \text { with leading exponent } a_{l_{1}}\right\} \\
\mathcal{S}_{l_{1}}^{\alpha_{1}}(h) & :=\left\{q \in \mathcal{S}_{l_{1}}(h) ; q \text { has leading term } c_{l_{1}}^{\alpha_{1}} x^{a_{l_{1}}}\right\} \\
m_{l_{1}}(h) & :=\text { number of roots in } \mathcal{S}_{l_{1}}(h), \text { counted with their multiplicities, } \\
m_{l_{1}}^{\alpha_{1}}(h) & :=\text { number of roots in } \mathcal{S}_{l_{1}}^{\alpha_{1}}(h), \text { counted with their multiplicities. }
\end{aligned}
$$

Then,

$$
\mathcal{S}_{l_{1}}(h)=\bigcup_{\alpha_{1}} \mathcal{S}_{l_{1}}^{\alpha_{1}}(h) \quad \text { and } \quad m_{l_{1}}(h)=\sum_{\alpha_{1}} m_{l_{1}}^{\alpha_{1}}(h)
$$

For simplicity, we shall denote

$$
m_{l_{1}}^{\alpha_{1}}(f)=m_{l_{1}}^{\alpha_{1}}, \quad m_{l_{1}}^{\alpha_{1}}(g)=n_{l_{1}}^{\alpha_{1}}, \quad \mathcal{S}_{l_{1}}(f) \cup \mathcal{S}_{l_{1}}(g)=\mathcal{S}_{l_{1}}
$$

Thus, for every $q \in \mathcal{S}_{l_{1}}^{\alpha_{1}}$ such that $q \not \equiv q_{1}$, there exist $a_{l_{1} l_{2}}^{\alpha_{1}}>a_{l_{1}}$ and $c_{l_{1} l_{2}}^{\alpha_{1} \alpha_{2}} \neq 0$ such that

$$
q(x)-q_{1}(x)=c_{l_{1} l_{2}}^{\alpha_{1} \alpha_{2}} x^{a_{l_{1} l_{2}}^{\alpha_{1}}}+\text { higher order terms. }
$$

We assume that the $a_{l_{1} l_{2}}^{\alpha_{1}}$ 's have been arranged in increasing order of magnitude, i.e.

$$
\begin{equation*}
a_{l_{1}}=a_{l_{1} 0}<a_{l_{1} 1}^{\alpha_{1}}<a_{l_{1} 2}^{\alpha_{1}}<\cdots<a_{l_{1} l_{2}}^{\alpha_{1}}<a_{l_{1}\left(l_{2}+1\right)}^{\alpha_{1}}<\cdots<\infty \tag{3.1}
\end{equation*}
$$

Define

$$
\mathcal{S}_{l_{1} l_{2}}^{\alpha_{1}}\left(h ; q_{1}\right):=\left\{q \in \mathcal{S}_{l_{1}}^{\alpha_{1}}\left(h ; q_{1}\right) ;\left(q-q_{1}\right) \text { has leading exponent } a_{l_{1} l_{2}}^{\alpha_{1}}\right\}
$$

and

$$
m_{l_{1} l_{2}}^{\alpha_{1}}\left(h ; q_{1}\right):=\text { number of roots in } \mathcal{S}_{l_{1} l_{2}}^{\alpha_{1}}\left(h ; q_{1}\right), \text { counted with their multiplicities. }
$$

As before, set

$$
m_{l_{1} l_{2}}^{\alpha_{1}}(f)=m_{l_{1} l_{2}}^{\alpha_{1}}, \quad m_{l_{1} l_{2}}^{\alpha_{1}}(g)=n_{l_{1} l_{2}}^{\alpha_{1}}, \quad \mathcal{S}_{l_{1} l_{2}}^{\alpha_{1}}\left(f ; q_{1}\right) \cup \mathcal{S}_{l_{1} l_{2}}^{\alpha_{1}}\left(g ; q_{1}\right)=\mathcal{S}_{l_{1} l_{2}}^{\alpha_{1}}
$$

Then,

$$
\sum_{l_{2} \geqslant 1} m_{l_{1} l_{2}}^{\alpha_{1}}\left(h ; q_{1}\right)=m_{l_{1}}^{\alpha_{1}}\left(h ; q_{1}\right)-M_{1}(h)
$$

where $M_{1}(h)$ is the multiplicity of $q_{1}$ as a root of $h$. In order to examine the integral in a neighborhood of $q_{1}$, we set

$$
\left|y-q_{1}(x)\right| \sim 2^{-p}, \quad p \geqslant a_{l_{1}} j
$$

Let us observe that this is equivalent to making the change of variables $x \mapsto x$ and $y \mapsto y-q_{1}(x)$. We denote this change of variables by $\varphi$. $\varphi$ will be a member of the class of admissible coordinates $\mathcal{C}_{0}$. Note that $\varphi$ preserves the origin and has Jacobian identically equal to 1 . In the subsequent analysis, the role of $k$, as depicted in cases 1 and 3 , is taken over by $p$. For example, we shall now need to consider the following ranges of $p$ and $j$ :

$$
a_{l_{1} l_{2}}^{\alpha_{1}} j \ll p \ll a_{l_{1}\left(l_{2}+1\right)}^{\alpha_{1}} j, \quad p \approx a_{l_{1} l_{2}}^{\alpha_{1}} j
$$

Subcase 1: $a_{l_{1} l_{2}}^{\alpha_{1}} j \ll p \ll a_{l_{1}\left(l_{2}+1\right)}^{\alpha_{1}} j$. In this region, the factors of $f$ and $g$ satisfy the following estimates:

$$
|y-q(x)|=\left|\left(y-q_{1}(x)\right)-\left(q(x)-q_{1}(x)\right)\right| \sim \begin{cases}2^{-a_{l_{1}^{\prime}} j}, & \text { if } q \in \mathcal{S}_{l_{1}^{\prime}}, l_{1}^{\prime}<l_{1} \\ 2^{-a_{l_{1}} j}, & \text { if } q \in \mathcal{S}_{l_{1}^{\prime}}, l_{1}^{\prime}>l_{1} \\ 2^{-a_{l_{1}} j}, & \text { if } q \in \mathcal{S}_{l_{1}} \backslash \mathcal{S}_{l_{1}}^{\alpha_{1}} \\ 2^{-a_{l_{1} l_{2}^{\prime}}^{\alpha_{1}} j}, & \text { if } q \in \mathcal{S}_{l_{1} l_{1}}^{\alpha_{1}}, l_{2}^{\prime} \leqslant l_{2} \\ 2^{-p}, & \text { if } q \in \mathcal{S}_{l_{1}}^{\alpha_{1} l_{2}^{\prime}}, l_{2}^{\prime}>l_{2} \\ 2^{-p}, & \text { if } q=q_{1}\end{cases}
$$

Therefore,

$$
|f(x, y)| \sim 2^{-\left(A_{l_{1}-1}+a_{l_{1}} B_{l_{1}}\right) j} 2^{-a_{l_{1}}\left(m_{l_{1}}-m_{l_{1}}^{\alpha_{1}}\right) j} 2^{-A_{l_{1} l_{2}}^{\alpha_{1}} j} 2^{-p B_{l_{1} l_{2}}^{\alpha_{1}}}
$$

where

$$
A_{l_{1} l_{2}}^{\alpha_{1}}=\sum_{l_{2}^{\prime} \leqslant l_{2}} a_{l_{1} l_{2}^{\prime}}^{\alpha_{1}} m_{l_{1} l_{2}^{\prime}}^{\alpha_{1}}, \quad B_{l_{1} l_{2}}^{\alpha_{1}}=\sum_{l_{2}^{\prime}>l_{2}} m_{l_{1} l_{2}^{\prime}}^{\alpha_{1}}
$$

i.e.

$$
\begin{aligned}
&|f(x, y)|\left.\sim 2^{-\left(A_{l_{1}-1}+a_{l_{1}} B_{l_{1}-1}-a_{l_{1}} m_{l_{1}}\right.}+A_{l_{1} l_{2}}^{\alpha_{1}}\right) j \\
& 2^{-B_{l_{1} l_{2}} p} \\
& \mid g(x, y) \sim 2^{-\left(C_{l_{1}-1}+a_{l_{1}} D_{l_{1}-1}-a_{l_{1}} n_{l_{1}}^{\alpha_{1}}+C_{l_{1} l_{2}}^{\alpha_{1}}\right) j} 2^{-D_{l_{1} l_{2}}^{\alpha_{1}} p} .
\end{aligned}
$$

This implies that,

$$
\begin{aligned}
& \sum_{(j, k, p) ; k \sim a_{l_{1}} j} \int \frac{|g(x, y)|^{\epsilon}}{|f(x, y)|^{\delta}} \chi_{j}(x) \chi_{k}(y) \chi_{p}\left(y-q_{1}(x)\right) d y d x \\
& a_{l_{1} l_{2}}^{\alpha_{1}} j \ll p \ll a_{l_{1}\left(l_{2}+1\right)}^{\alpha_{1}} j \\
& \sim \sum_{\substack{(j, k, p) ; k \sim a_{1} j \\
a_{l_{1}}^{\alpha_{1}} j \ll p \ll a_{l_{1}\left(l_{2}+1\right)}^{\alpha_{1}} j}} \frac{2^{-\left(C_{l_{1}-1}+a_{l_{1}} D_{l_{1}-1}-a_{l_{1}} n_{l_{1}}^{\alpha_{1}}+C_{l_{1} l_{2}}^{\alpha_{1}}\right) j \epsilon} 2^{-D_{l_{1} l_{2}}^{\alpha_{1}} p \epsilon}}{2^{-\left(A_{l_{1}-1}+a_{l_{1}} B{l_{1}-1}^{\alpha_{1}} a_{l_{1}} m_{l_{1}}^{\alpha_{1}}+A_{l_{1} l_{2}}^{\alpha_{1}}\right) j \delta} 2^{-B_{l_{1} l_{2}}^{\alpha_{1}} p \delta}} 2^{-p} 2^{-j} \\
& \sim \sum_{\substack{(j, p) ; \\
a_{l_{1} l_{2}}^{\alpha_{1}} j \ll p \ll a_{l_{1}\left(l_{2}+1\right)}^{\alpha_{1}}}} 2^{j \delta\left(A_{l_{1}-1}+a_{l_{1}} B_{l_{1}-1}-a_{l_{1}} m_{l_{1}}^{\alpha_{1}}+A_{l_{1} l_{2}}^{\alpha_{1}}\right)} \\
& \times 2^{-j\left[\epsilon\left(C_{l_{1}-1}+a_{l_{1}} D_{l_{1}-1}-a_{l_{1}} n_{l_{1}}^{\alpha_{1}}+C_{l_{1} l_{2}}^{\alpha_{1}}\right)+1\right]} 2^{p\left(B_{l_{1} l_{2}}^{\alpha_{1}} \delta-D_{l_{1} l_{2}}^{\alpha_{1}}-1\right)} .
\end{aligned}
$$

Just as before, we can show that the infinite series converges:
(a) for $\delta B_{l_{1} l_{2}}^{\alpha_{1}}-\epsilon D_{l_{1} l_{2}}^{\alpha_{1}}-1 \leqslant 0$ iff

$$
\begin{array}{r}
\delta<\frac{\epsilon\left(C_{l_{1}-1}+a_{l_{1}} D_{l_{1}-1}-a_{l_{1}} n_{l_{1}}^{\alpha_{1}}+C_{l_{1} l_{2}}^{\alpha_{1}}+a_{l_{1} l_{2}}^{\alpha_{1}} D_{l_{1} l_{2}}^{\alpha_{1}}\right)+a_{l_{1} l_{2}}^{\alpha_{1}}+1}{A_{l_{1}-1}+a_{l_{1}} B_{l_{1}-1}-a_{l_{1}} m_{l_{1}}^{\alpha_{1}}+A_{l_{1} l_{2}}^{\alpha_{1}}+a_{l_{1} l_{2}}^{\alpha_{1}} B_{l_{1} l_{2}}^{\alpha_{1}}} \\
:=\omega\left(q_{1} ; l_{1}, l_{2} ; \alpha_{1}\right)
\end{array}
$$

(b) for $\delta B_{l_{1} l_{2}}^{\alpha_{1}}-\epsilon D_{l_{1} l_{2}}^{\alpha_{1}}-1>0$ iff

$$
\begin{array}{r}
\delta<\frac{\epsilon\left(C_{l_{1}-1}+a_{l_{1}} D_{l_{1}-1}-a_{l_{1}} n_{l_{1}}^{\alpha_{1}}+C_{l_{1} l_{2}}^{\alpha_{1}}+a_{l_{1} l_{2}}^{\alpha_{1}} D_{l_{1}\left(l_{2}+1\right)}^{\alpha_{1}}\right)+a_{l_{1}\left(l_{2}+1\right)}^{\alpha_{1}}+1}{A_{l_{1}-1}+a_{l_{1}} B_{l_{1}-1}-a_{l_{1}} m_{l_{1}}^{\alpha_{1}}+A_{l_{1}\left(l_{2}+1\right)}^{\alpha_{1}}+a_{l_{1}\left(l_{2}+1\right)}^{\alpha_{1}} B_{l_{1}\left(l_{2}+1\right)}^{\alpha_{1}}} \\
:=\omega\left(q_{1} ; l_{1},\left(l_{2}+1\right) ; \alpha_{1}\right) .
\end{array}
$$

We observe that

$$
\left(\frac{a_{l_{1} l_{2}}^{\alpha_{1}}+1}{A_{l_{1}-1}+a_{l_{1}} B_{l_{1}-1}-a_{l_{1}} m_{l_{1}}^{\alpha_{1}}+A_{l_{1} l_{2}}^{\alpha_{1}}+a_{l_{1} l_{2}}^{\alpha_{1}} B_{l_{1} l_{2}}^{\alpha_{1}}}\right)^{-1}:=\left(\delta_{l_{1} l_{2}}^{\alpha_{1}}\right)^{-1}
$$

is the point of intersection of the bisectrix $p=q$ with the face of the generalized Newton diagram of the function $\tilde{f}(x, y)=f \circ \varphi(x, y)=f\left(x, y-q_{1}(x)\right)$ that has slope $\frac{-1}{a_{l_{1} l_{2}}^{\alpha_{1}}}$. Denoting by $\left(\tilde{\delta}_{l_{1} l_{2}}^{\alpha_{1}}\right)^{-1}$ the corresponding quantity for $g$, we get

$$
\omega\left(q_{1} ; l_{1}, l_{2} ; \alpha_{1}\right)=\delta_{l_{1} l_{2}}^{\alpha_{1}}\left(1+\frac{\epsilon}{\tilde{\delta}_{l_{1} l_{2}}^{\alpha_{1}}}\right) .
$$

Now, using the same techniques as in case 2 , we see that the optimal exponent $\delta$ satisfies

$$
\delta<\min _{l} \delta_{l}(\varphi)\left(1+\frac{\epsilon}{\tilde{\delta}_{l}(\varphi)}\right)
$$

Subcase 2: $p \gg a_{l_{1} L_{2}}^{\alpha_{1}} j$, where $L_{2}=\max l_{2}$. This is the step that deals with the region "close" to the root $q_{1}$ and highlights the cancellation (if any) of the factor
( $y-q_{1}(x)$ ) between $|g|^{\epsilon}$ and $|f|^{\delta}$. In this region, $f$ and $g$ satisfy the following size estimates:

$$
\begin{aligned}
|f(x, y)| & \sim 2^{-\left(A_{l_{1}-1}+a_{l_{1}} B_{l_{1}}\right) j} 2^{-a_{l_{1}}\left(m_{l_{1}}-m_{l_{1}}^{\alpha_{1}}\right) j} 2^{-A_{l_{1} L_{2}}^{\alpha_{1}} j} 2^{-p M_{1}(f)}, \\
|g(x, y)| & \sim 2^{-\left(C_{l_{1}-1}+a_{l_{1}} D_{l_{1}}\right) j} 2^{-a_{l_{1}}\left(n_{l_{1}}-n_{l_{1}}^{\alpha_{1}}\right) j} 2^{-C_{l_{1} L_{2}}^{\alpha_{1}} j} 2^{-p M_{1}(g)}
\end{aligned}
$$

So routine computations yield

$$
\begin{array}{r}
\sum_{\substack{(j, k, p) ; k \sim a_{l} j \\
p \gg a_{l_{1} L_{2}}^{\alpha_{1}} j}} \int \frac{|g(x, y)|^{\epsilon}}{|f(x, y)|^{\delta}} \chi_{j}(x) \chi_{k}(y) \chi_{p}\left(y-q_{1}(x)\right) d y d x \\
\sim \sum_{\substack{(j, k, p) ; k \sim a_{l} j \\
p \gg a_{l_{1} L_{2}}^{\alpha_{1}} j}} 2^{p\left(\delta M_{1}(f)-\epsilon M_{1}(g)-1\right)} 2^{j \delta\left(A_{l_{1}-1}+a_{l_{1}} B_{l_{1}-1}-a_{l_{1}} m_{l_{1}}^{\alpha_{1}}+A_{l_{1} L_{2}}^{\alpha_{1}}\right)} \\
\end{array}
$$

which converges iff

$$
\delta<\min \left(\frac{\epsilon M_{1}(g)+1}{M_{1}(f)}, \omega\left(q_{1} ; l_{1}, L_{2} ; \alpha_{1}\right)\right)
$$

We observe that $\left(\epsilon M_{1}(g)+1\right) / M_{1}(f)$ is the value of $\omega$ corresponding to the horizontal faces of the generalized Newton diagrams of $\tilde{f}$ and $\tilde{g}$.

Subcase 3: $p \approx a_{l_{1} l_{2}}^{\alpha_{1}} j$. This range of indices depicts the region near the cluster of roots $\mathcal{S}_{l_{1} l_{2}}^{\alpha_{1}}$. Let us fix a real root $q_{2}$ in this cluster. Then,

$$
q_{2}(x)-q_{1}(x)=c_{l_{1} l_{2}}^{\alpha_{1} \alpha_{2}} x^{a_{l_{1} l_{2}}^{\alpha_{1}}}+\text { higher order terms }
$$

for some index $\alpha_{2}$ such that $c_{l_{1} l_{2}}^{\alpha_{1} \alpha_{2}} \neq 0$.
Defining
$\mathcal{S}_{l_{1} l_{2}}^{\alpha_{1} \alpha_{2}}\left(h ; q_{1}\right):=\left\{q \in \mathcal{S}_{l_{1} l_{2}}^{\alpha_{1}}\left(h ; q_{1}\right) ;\left(q-q_{1}\right)(x)=c_{l_{1} l_{2}}^{\alpha_{1} \alpha_{2}} x^{a_{l_{1} l_{2}}^{\alpha_{1}}}+\right.$ higher order terms $\}$, and setting

$$
\mathcal{S}_{l_{1} l_{2}}^{\alpha_{1} \alpha_{2}}:=\mathcal{S}_{l_{1} l_{2}}^{\alpha_{1} \alpha_{2}}\left(f ; q_{1}\right) \cup \mathcal{S}_{l_{1} l_{2}}^{\alpha_{1} \alpha_{2}}\left(g ; q_{1}\right),
$$

we see that this subcase is identical to case 3 , except that we now view the original cluster of roots under a finer resolution; namely, we consider only those roots which have the same leading term as $q_{1}$ and whose difference from $q_{1}$ has a leading factor of $c_{l_{1} l_{2}}^{\alpha_{1} \alpha_{2}} x^{a_{l_{1} l_{2}}^{\alpha_{1}}}$.

Thus, as before, we have that for any $q \in \mathcal{S}_{l_{1} l_{2}}^{\alpha_{1} \alpha_{2}}$ such that $q \not \equiv q_{2}$, there exist $a_{l_{1} l_{2} l_{3}}^{\alpha_{1} \alpha_{2}}>a_{l_{1} l_{2}}^{\alpha_{1}}$ and $c_{l_{1} l_{2} l_{3}}^{\alpha_{1} \alpha_{3} \alpha_{3}} \neq 0$ such that

$$
\left(q-q_{2}\right)(x)=c_{l_{1} l_{2} l_{3}}^{\alpha_{1} \alpha_{2} \alpha_{3}} x^{a_{l_{1} l_{2} l_{3}}^{\alpha_{1} \alpha_{2}}}+\text { higher order terms }
$$

Let us order

$$
a_{l_{1} l_{2}}^{\alpha_{1}}<a_{l_{1} l_{2} 1}^{\alpha_{1} \alpha_{2}}<\cdots<a_{l_{1} l_{2} l_{3}}^{\alpha_{1} \alpha_{2}}<a_{l_{1} l_{2}\left(l_{3}+1\right)}^{\alpha_{1} \alpha_{2}}<\cdots<\infty .
$$

Now, setting $y-q_{2}(x) \sim 2^{-p_{2}}, p_{2} \geqslant p$ (in other words, making the transformation $\left.x \mapsto x, y \mapsto y-q_{2}(x)\right)$, we may treat this subcase exactly as the parent case 3 by further subdividing into cases depending on the range of $p_{2}$, i.e.

$$
a_{l_{1} l_{2} l_{3}}^{\alpha_{1} \alpha_{2}} j \ll p_{2} \ll a_{l_{1} l_{2}\left(l_{3}+1\right)}^{\alpha_{1} \alpha_{2}} j \quad \text { and } \quad p_{2} \sim a_{l_{1} l_{2} l_{3}}^{\alpha_{1} \alpha_{2}} j .
$$

The second type of region gives rise to further subcases of its own (identical to the ones already discussed but with an added degree of refinement in the cluster). In general, one considers, following an inductive scheme, progressively finer clusters given by

$$
\begin{align*}
& \mathcal{S}_{l_{1} l_{2} \ldots l_{s}}^{\alpha_{1} \alpha_{2} \ldots \alpha_{s-1}}  \tag{3.2}\\
& \quad:=\mathcal{S}_{l_{1} l_{2} \ldots l_{s}}^{\alpha_{1} \alpha_{2} \ldots \alpha_{s-1}}\left(f ; q_{1}, q_{2}, \cdots, q_{s-1}\right) \cup \mathcal{S}_{l_{1} l_{2} \ldots l_{s}}^{\alpha_{1} \alpha_{2} \ldots \alpha_{s-1}}\left(g ; q_{1}, q_{2}, \cdots, q_{s-1}\right)
\end{align*}
$$

and

$$
\begin{equation*}
\mathcal{S}_{l_{1} l_{2} \ldots l_{s}}^{\alpha_{1} \alpha_{2} \ldots \alpha_{s}}:=\mathcal{S}_{l_{1} l_{2} \ldots l_{s}}^{\alpha_{1} \alpha_{2} \ldots \alpha_{s}}\left(f ; q_{1}, q_{2}, \cdots, q_{s-1}\right) \cup \mathcal{S}_{l_{1} l_{2} \ldots l_{s}}^{\alpha_{1} \alpha_{2} \ldots \alpha_{s}}\left(g ; q_{1}, q_{2}, \cdots, q_{s-1}\right), \tag{3.3}
\end{equation*}
$$

where the sets in (3.2) and (3.3) are defined recursively as follows.
Given $\mathcal{S}_{l_{1} \cdots l_{s-1}}^{\alpha_{1} \cdots \alpha_{s-1}}\left(h ; q_{1}, \cdots, q_{s-2}\right)$, we set

$$
\begin{aligned}
& \mathcal{S}_{l_{1} l_{2} \cdots l_{s}}^{\alpha_{1} \alpha_{2} \cdots \alpha_{s-1}}\left(h ; q_{1}, q_{2}, \ldots, q_{s-1}\right):=\left\{q \in \mathcal{S}_{l_{1} l_{2} \cdots l_{s-1}}^{\alpha_{1} \alpha_{2} \cdots \alpha_{s-1}}\left(h ; q_{1}, q_{2}, \cdots, q_{s-2}\right)\right. \\
&\left.\left(q-q_{s-1}\right)(x) \text { has leading exponent } a_{l_{1} l_{2} \cdots l_{s}}^{\alpha_{1} \alpha_{2} \cdots \alpha_{s-1}}\right\}
\end{aligned}
$$

and

$$
\begin{aligned}
& \mathcal{S}_{l_{1} l_{2} \cdots l_{s}}^{\alpha_{1} \alpha_{2} \cdots \alpha_{s}}\left(h ; q_{1}, \cdots, q_{s-1}\right):=\left\{q \in \mathcal{S}_{l_{1} l_{2} \cdots l_{s}}^{\alpha_{1} \alpha_{2} \cdots \alpha_{s-1}}\left(h ; q_{1}, \cdots, q_{s-1}\right)\right. \\
& \left.\quad\left(q-q_{s-1}\right)(x)=c_{l_{1} l_{2} \cdots l_{s}}^{\alpha_{1} \alpha_{2} \cdots \alpha_{s}} x^{a_{l_{1} l_{2} \cdots l_{s}}^{\alpha_{1} \alpha_{2} \cdots \alpha_{s-1}}}+\text { higher order terms, } c_{l_{1} l_{2} \cdots l_{s}}^{\alpha_{1} \alpha_{2} \cdots \alpha_{s}} \neq 0\right\}
\end{aligned}
$$

Here $h=f$ or $g$, and $q_{i}$ is a fixed real root of the cluster $\mathcal{S}_{l_{1} l_{2} \cdots l_{i}}^{\alpha_{1} \alpha_{2} \cdots \alpha_{i}}$, for $1 \leqslant i \leqslant s-1$. One observes that

$$
\mathcal{S}_{l_{1}} \supseteq \mathcal{S}_{l_{1}}^{\alpha_{1}} \supsetneq \mathcal{S}_{l_{1} l_{2}}^{\alpha_{1}} \supseteq \mathcal{S}_{l_{1} l_{2}}^{\alpha_{1} \alpha_{2}} \supsetneq \cdots
$$

so the process terminates in a finite number of steps, when each of the final clusters $\mathcal{S}_{l_{1} l_{2} \cdots l_{s}}^{\alpha_{1} \alpha_{2} \cdots \alpha_{s}}$ consists of roots that are identical, i.e. when the degree of resolution is high enough to distinguish between any two distinct roots.

Case 4: Complex roots. We shall now indicate how the above arguments are modified in the presence of complex roots. Let a factor of $f$ or $g$ be given by $y-q_{0}(x)$, and let

$$
q_{0}(x)=c x^{a_{l}}+\text { higher powers }
$$

Here $l=l_{1}$ and $c=c_{l_{1}}^{\alpha_{1}}$ for some $\alpha_{1}$. There are two possibilities: $c \in \mathbb{R}$ and $c \notin \mathbb{R}$.
For $c \in \mathbb{R}$,

$$
\begin{aligned}
\left|y-q_{0}(x)\right| & \sim\left|y-\Re q_{0}(x)\right|+\left|\Im q_{0}(x)\right| \\
& \sim\left|y-\Re q_{0}(x)\right|+|x|^{b_{l_{1}}} \quad \text { for some } b_{l_{1}}>a_{l_{1}}
\end{aligned}
$$

Now, in the ranges $k \ll a_{l_{1}} j$ and $k \gg a_{l_{1}} j$, the bounds for $f$ and $g$ from above and below are unaffected by the presence of $|x|^{b_{l_{1}}}$, since there is no cancellation between $y$ and $|x|^{a_{l_{1}}}$, and $|x|^{a_{l_{1}}}$ is much bigger than $|x|^{b_{l_{1}}}$. Thus we need only consider the range $k \sim a_{l_{1}} j$. We set

$$
y-\Re q_{0}(x) \sim 2^{-m}, \quad m \geqslant a_{l_{1}} j
$$

Then,

$$
\left|y-q_{0}(x)\right| \sim 2^{-\min \left(b_{l_{1}} j, m\right)}
$$

Let us consider the collection of leading exponents of $q(x)-\Re q_{0}(x)$, where $q \in \mathcal{S}_{l_{1}}^{\alpha_{1}}$, and order them as in (3.1). We know that $b_{l_{1}}$ occurs in this sequence, as it is the leading exponent of $q_{0}(x)-\Re q_{0}(x)$, so let us assume

$$
b_{l_{1}}=a_{l_{1} l_{2}}^{\alpha_{1}} \quad \text { for some } l_{2} .
$$

Now, depending on the relative sizes of $m$ and $b_{l_{1}} j$, we consider subcases of the following forms:

It is a matter of direct calculation to verify that the size estimates of $f$ and $g$ that one obtains in cases of the form (1) and (3) are identical to the ones obtained in subcase 1 of case 3 , while the treatment of cases (2) and (4) exactly matches that of subcase 3 of case 3 . Thus, a repetition of the same computations show that the optimal $\delta$ satisfies the following inequality:

$$
\delta<\min _{l}\left[\delta_{l}(\tilde{\varphi})\left(1+\frac{\epsilon}{\tilde{\delta}_{l}(\tilde{\varphi})}\right)\right]
$$

where $\tilde{\varphi}$ is an admissible change of coordinates given by $(x, y) \mapsto\left(x, y-\Re q_{0}(x)\right)$.
This concludes our treatment of the case when $c \in \mathbb{R}$. The other case is in fact simpler than the preceding one, because now we have

$$
\left|y-q_{0}(x)\right| \sim\left(|y|+|x|^{a_{l_{1}}}\right)
$$

The size estimates for $f$ and $g$ are, of course, unaffected in the range $a_{l_{1}} j \ll k \ll$ $a_{l_{1}+1} j$. Even in the range $k \sim a_{l_{1}} j$, there can be no cancellation between the terms $|y|$ and $|x|^{a_{l_{1}}}$, so this range does not require a separate treatment. The case of complex roots has therefore been verified.

It is now possible to list all the members of $\mathcal{C}_{0}$, as follows:
$\mathcal{C}_{0}:=\left\{\begin{array}{l}\varphi:(x, y) \mapsto\left(x^{\prime}, y^{\prime}\right) ; \varphi \text { may be of the following forms: } \\ \text { either }\left\{\begin{array}{l}x^{\prime}=x \\ y^{\prime}=y-q(x)\end{array} \quad \text { or } \quad \begin{cases}x^{\prime}=x-q(y) \\ y^{\prime} & =y\end{cases} \right. \\ \left.\begin{array}{l}\text { where } q \text { is a real-valued Puiseux series with leading exponent } \geqslant 1 . \\ y-q(x) \text { or } x-q(y) \text { may be either a real root of } f \text { or } g, \text { or the } \\ \text { real part of a complex root of } f \text { or } g, \text { provided the complex root } \\ \text { has a real leading coefficient. }\end{array}\right\} .\end{array}\right.$
The proof shows that the $\delta$ given by

$$
\delta_{0}^{\mathcal{C}_{0}}(g, f, \epsilon):=\min _{\varphi \in \mathcal{C}_{0}} \delta_{0}(g, f, \epsilon ; \varphi)
$$

is, in fact, optimal. Moreover, it is not hard to see that the above analysis may be replicated even if we set

$$
y-t(x) \sim 2^{-p}
$$

where $t$ is a convergent Puiseux series, not necessarily associated to $f$ or $g$. This shows that, in fact,

$$
\delta_{0}^{\mathcal{C}}(g, f, \epsilon)=\delta_{0}^{\mathcal{C}_{0}}(g, f, \epsilon)
$$

To complete the proof of Theorem we need to justify the assumption $a_{l} \geqslant 1$. This is clearly the case in the range $k \geqslant j-K$ for a large enough constant $K$. In the remaining range $k \leqslant j-K$ it suffices to write $f$ and $g$ as polynomials in $x$ with analytic coefficients in $y$ (up to the usual non-vanishing factors). The zeroes of $f$ and $g$ are then of the form $x=\tilde{r}(y)$, where $\tilde{r}(y)$ is a Puiseux series in $y$ with leading exponent greater than or equal to 1 . We can now repeat our arguments with the roles of $x$ and $y$ interchanged. The proof of Theorem 1 is therefore complete.

Remark. It is interesting to note the contrast of the weighted situation with the unweighted case, where the transformation yielding the optimal $\delta$ always turned out to be analytic [2]. For the weighted integral, the change of coordinates required to achieve the optimal $\delta$, while being algebraic and indeed of a very special form, need not be better than once continuously differentiable, as the following example shows:

Let us consider the integral

$$
I=\iint_{B} \frac{\left|\left(y^{3}-x^{4}\right)\left(y-x^{2}\right)^{3}\right|}{\left|\left(y^{3}-x^{4}\right)\left(y^{3}-x^{4}-x^{5}\right)\left(y-x^{2}\right)\right|^{\delta}} d y d x
$$

It is easy to see that the optimal $\delta$ is 1 and is obtained using the transformation $(x, y) \mapsto\left(x, y-x^{\frac{4}{3}}(1+x)^{\frac{1}{3}}\right)$. Moreover, no analytic transformation of the form $(x, y) \mapsto(x, y-s(x))$ works, as can be seen by setting $s(x)=c_{1} x^{p}+c_{2} x^{q}+\cdots$, where $q>p \geqslant 1$ are integers, and computing directly the quantity $\min _{l} \delta_{l}\left(1+\tilde{\delta}_{l}^{-1}\right)$. This involves having to consider various ranges of values of $c_{1}, c_{2}$ and $p$, whence one obtains that

$$
\min _{l} \delta_{l}\left(1+\tilde{\delta}_{l}^{-1}\right)=\left\{\begin{array}{ll}
\frac{8}{7} & \text { if } p=1 \\
\min \left(\frac{31}{28}, \frac{4 q+5}{q+8}\right) & \text { if } p=2, c_{1}=1, c_{2} \neq 0 ; \\
\frac{31}{28} & \text { otherwise }
\end{array}\right\}
$$

all of which are larger than 1 .
Example. Finally, we would like to analyze Varchenko's counterexample, which was one of the motivating factors for the weighted integral, in this set-up. When $\lambda$ is non-negative, the only admissible transformation for the integral in (1.3) is the identity transformation, as the denominator has neither a real root, nor a complex root with a real leading coefficient.

Thus, when $\lambda=0$, the Newton diagram of $f$ consists of a single boundary segment joining the points $(0,4)$ and $(8,0)$. One has

$$
\delta_{1}=\frac{3}{8}, \quad \tilde{\delta}_{1}=\frac{3}{2} .
$$

Therefore,

$$
\delta_{0}^{\mathcal{C}}(g, f, 1)=\frac{3}{8}\left(1+\frac{2}{3}\right)=\frac{5}{8} .
$$

When $\lambda>0$, the Newton diagram of $f$ again has a single boundary segment, joining $(0,4)$ and $(4,0)$, and one obtains

$$
\delta_{1}=\frac{1}{2}, \quad \tilde{\delta}_{1}=2
$$

This implies

$$
\delta_{0}^{\mathcal{C}}(g, f, 1)=\frac{1}{2}\left(1+\frac{1}{2}\right)=\frac{3}{4}
$$

When $\lambda<0$, there exist non-trivial admissible transformations, all of which are of the form $(x, y) \mapsto(x, y-q(x))$, for some Puiseux series $q$ with leading exponent 1. For any non-trivial $\varphi \in \mathcal{C}_{0}$, the generalized Newton diagram of $f \circ \varphi$ therefore has two boundary segments-one joining $(0,4)$ and $(2,2)$, and the other joining $(2,2)$ and $(4 p, 0)$. The generalized Newton diagram of $g \circ \varphi$ has a single segment, joining $(0,1)$ and $(1,0)$. The computations yield

$$
\delta_{1}(\varphi)=\frac{1}{2}, \quad \tilde{\delta}_{1}(\varphi)=2, \quad \delta_{2}(\varphi)=\frac{1}{2}, \quad \tilde{\delta}_{2}(\varphi)=2 p
$$

Therefore,

$$
\delta_{0}^{\mathcal{C}}(g, f, 1)=\min \left(\frac{1}{2}\left(1+\frac{1}{2}\right), \frac{1}{2}\left(1+\frac{1}{2 p}\right)\right)=\frac{1}{2}+\frac{1}{4 p} .
$$

While all of the above cases corroborate Varchenko's result, the case $\lambda<0$ deserves special attention. The above analysis provides the optimal $\delta$ even in this case, in contrast to [5], where it was only shown that

$$
\delta<\frac{2 p}{4 p-2}
$$

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