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## TRANSSERIES FOR BEGINNERS


#### Abstract

From the simplest point of view, transseries concern manipulations on formal series, or a new kind of expansion for real-valued functions. But transseries constitute much more than that - they have a very rich (algebraic, combinatorial, analytic) structure. The set of transseriesalso known as the transline - is a large ordered field, extending the real number field, and endowed with additional operations such as exponential, logarithm, derivative, integral, composition. Over the course of the last 20 years or so, transseries have emerged in several areas of mathematics: dynamical systems, model theory, computer algebra, surreal numbers. This paper is an exposition for the non-specialist mathematician.


All a mathematician needs to know in order to apply transseries.

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## Introduction.

Although transseries expansions are prominent in certain areas of mathematics, they are not well known to mathematicians in general. Here, I try to bring these beautiful mathematical structures to the attention of non-specialists. This paper complements the already-existing survey articles such as [1, 49], or monographs [37, 43, 54].

Transseries come in various flavors. Here I focus on one particular variantthe real, grid-based transseries-since they are the ones which are most amenable to explicit computations, and transseries representing real-valued functions naturally arising in analysis (e.g., as solutions to algebraic differential equations) are usually of this type. Once familiar with one variant, it should be relatively easy to work with another.

The major part of this paper (Section 3) presents a formal construction of the differential field of real, grid-based transseries. Section 4 illustrates its use in practice through examples: transseries expansions for functional inverses, for anti-derivatives, for solutions of differential equations, etc. The development is entirely formal; the analytic aspects and origins of the subject (computer algebra limit algorithms, Écalle's generalization of Borel summation, Hardy fields, etc.) are omitted-a survey of that aspect of the subject would warrant a separate paper. This restriction allows for a self-contained exposition, suited for mathematicians regardless of their specialties.

There are several constructions of the various fields of transseries already in the literature, smoothing out and filling in details in the original papers; for example: van den Dries-Macintyre-Marker [18], van der Hoeven [37], and Costin [11]. These all require a certain technical apparatus, despite the simplicity of the basic construction. Here I try to avoid such requirements and assume only a minimum of background knowledge.

Sections 1 and 2, which are intended to lure the reader into the transseries world, give examples of natural computations which can be made precise in this framework. Section 3 deals with the rigorous construction of grid-based transseries. It is intended as: All a mathematician needs to know in order to
apply transseries. Section 4 contains worked-out examples, partly computed with the aid of computer algebra software. Section 5 gives suggestions for further reading. (This introduction is taken mostly from an anonymous referee's report for an earlier draft of the paper. That referee understood what this paper is about better than I did myself!)

Review of "Fraktur" style letters:

$$
(A \mathfrak{A})(B \mathfrak{B})(G \mathfrak{G})(J \mathfrak{J})(M \mathfrak{M})(a \mathfrak{a})(b \mathfrak{b})(g \mathfrak{g})(l \mathfrak{l})(m \mathfrak{m})(n \mathfrak{n})
$$

## 1 Sales Pitch.

One day long ago, I wrote Stirling's formula like this:

$$
\log \Gamma(x)=\left(x-\frac{1}{2}\right) \log x-x+\frac{\log (2 \pi)}{2}+\sum_{n=1}^{\infty} \frac{B_{2 n}}{2 n(2 n-1) x^{2 n-1}}
$$

where the $B_{2 n}$ are the Bernoulli numbers. But my teacher gently told me that the series diverges for every $x$. What a disappointment!

Leonhard Euler [25, p. 220] (the master of us all [19]) wrote:

$$
\sum_{j=0}^{\infty} \frac{(-1)^{j} j!}{x^{j+1}}=-e^{x} \operatorname{Ei}(-x)
$$

where the exponential integral function is defined by

$$
\operatorname{Ei}(-x):=\int_{-\infty}^{-x} \frac{e^{t}}{t} d t
$$

But later mathematicians sneered at this, saying that the series diverges wildly.
To study a sequence $a_{j}$, it is sometimes useful to consider the "generating function"

$$
\sum_{j=0}^{\infty} a_{j} z^{j}=\sum_{j=0}^{\infty} \frac{a_{j}}{x^{j}}
$$

(The change of variables $z=1 / x$ was made so that we can consider not $z$ near zero but $x$ near infinity, as we will always do here.) In fact, it is quite useful to consider such a series "formally" even if the series diverges [60]. The generating function for the sequence $2^{j}$ is of course

$$
\sum_{j=0}^{\infty} \frac{2^{j}}{x^{j}}=\frac{1}{1-2 / x}
$$

But who among you has not secretly substituted $x=1$ to get

$$
\sum_{j=0}^{\infty} 2^{j}=-1
$$

and wondered at it?
To study asymptotic behavior of functions, G. H. Hardy [30] promoted the class $L$ of "orders of infinity": all functions (near $\infty$ ) obtained starting with constants and $x$, then applying the field operations, exp, and log repeatedly in any order. Function $x e^{x}$ is a valid member of that class. Liouville had shown that its inverse function isn't (see [6]). What cruel classification would admit a function but not its inverse?

Undergraduate courses in ordinary differential equations tell us how to solve a linear differential equation with analytic coefficients in terms of power series-at least at ordinary points, and at regular singular points. But power series solutions do not work at irregular singular points. Is it hopeless to understand solutions near these points?

Solving linear differential equations with constant coefficients can be approached by factoring of operators. Take, for example, $3 Y^{\prime \prime}-5 Y^{\prime}-2 Y=3 x$. Writing $\partial$ for the derivative operator and $I$ for the identity operator, this can be written $L[Y]=3 x$, where $L=3 \partial^{2}-5 \partial-2 I$. Then factor this polynomial, $L=3(\partial-2 I)(\partial+(1 / 3) I)$ and solve $L[Y]=3 x$ with two successive integrations: First write $Y_{1}=(\partial+(1 / 3) I) Y$. Then solve $\partial Y_{1}-2 Y_{1}=x$ to get $Y_{1}=A e^{2 x}-1 / 4-x / 2$. Then solve $\partial Y+(1 / 3) Y=A e^{2 x}-1 / 4-x / 2$ to get $Y=(3 A / 7) e^{2 x}+B e^{-x / 3}+15 / 4-3 x / 2$. Wouldn't it be grand if this could be done for linear differential equations with variable coefficients? But we cannot solve the differential equation $Y^{\prime \prime}+x Y^{\prime}+Y=0$ by factoring $\partial^{2}+x \partial+I=(\partial-\alpha(x) I)(\partial-\beta(x) I)$, where $\alpha(x)$ and $\beta(x)$ are polynomials; or rational functions; or elementary functions. (The "commutative" factoring method-the quadratic formula-does not work, since $\partial \beta(x) I \neq \beta(x) \partial$ if $\beta$ is not constant.) But what if we could factor with some new, improved, simple, versatile class of functions?

Well, brothers and sisters, I am here today to tell you: If you love these formulas, you need no longer hide in the shadows! The answer to all of these woes is here.

## Transseries

The differential field of transseries was discovered [or, some would say, invented] independently in various parts of mathematics: dynamical systems, model theory, computer algebra, surreal numbers. Some feel it was surprisingly recent for something so natural. Roots of the subject go back to Écalle
[21] and $\mathrm{Il}^{\prime}$ yashenko [40] working in dynamical systems; Dahn and Göring $[15,16]$ working in model theory; Geddes \& Gonnet [27] working in computer algebra; Kruskal working in surreal numbers (unpublished: see the Epilog in the Second Edition of On Numbers and Games [10]). They arrived at eerily similar mathematical structures, although they did not have all the same features. It is Ecalle who recognized the power of these objects, coined the term, developed them systematically and in their own right, found "the" way to associate functions to them. [I am not tracing the history here. Precursorsin addition to Hardy [30], Levi-Cevita [45], du Bois-Reymond [3], even Euler [25]-include Lightstone \& Robinson [46], Salvy \& Shackell [52], Rosenlicht [51], and Boshernitzan [4]. This listing is far from complete: Additional historical remarks are in $[37,43,49,55]$.]

I hope this paper will show that knowledge of model theory or dynamical systems or computer algebra or surreal numbers is not required in order to understand this new, beautiful, complex object.

In this paper, we consider only series used for $x \rightarrow+\infty$. Limits at other locations, and from other directions, are related to this by a change of variable. For example, to consider $z \rightarrow 1$ from the left, write $z=1-1 / x$ or $x=1 /(1-z)$.

## 2 What Is a Transseries?

There is an ordered group $\mathfrak{G}$ of transmonomials and a differential field $\mathbb{T}$ of transseries. But $\mathfrak{G}$ and $\mathbb{T}$ are each defined in terms of the other, in the way logicians like to do. There is even some spiffy notation (taken from [37]): $\mathbb{T}=\mathbb{R} \llbracket \mathfrak{G} \rrbracket=\mathbb{R} \llbracket x \rrbracket$. The definition is carried out formally in Section 3. But for now let's see informally what they look like. [This is "informal" since, for example, some terms are used before they are defined, so that the whole thing is circular.]
(a) A log-free transmonomial has the form $x^{b} e^{L}$, where $b$ is real and $L$ is a purely large log-free transseries; " $x$ " and " $e$ " are just symbols. Examples:

$$
x^{-1}, \quad x^{\pi} e^{x^{\sqrt{2}}-3 x}, \quad e^{\sum_{j=0}^{\infty} x^{-j} e^{x}}
$$

Use $x^{b_{1}} e^{L_{1}} \cdot x^{b_{2}} e^{L_{2}}=x^{b_{1}+b_{2}} e^{L_{1}+L_{2}}$ for the group operation "multiplication" and group identity $x^{0} e^{0}=1$. The ordering $\succ$ (read "far larger than", sometimes written $\gg$ instead) is defined for $\mathfrak{G}$ lexicographically: $x^{b_{1}} e^{L_{1}} \succ x^{b_{2}} e^{L_{2}}$ iff $L_{1}>L_{2}$ or $\left\{L_{1}=L_{2}\right.$ and $\left.b_{1}>b_{2}\right\}$. As usual, the converse relation $\prec$ is called "far smaller than". Examples:

$$
e^{\sum_{j=0}^{\infty} x^{-j} e^{x}} \succ e^{x} \succ x^{-3} e^{x} \succ x^{\pi} \succ x^{-1} \succ x^{-5} \succ x^{2010} e^{-x}
$$

(b) A log-free transseries is a (possibly infinite) formal sum $T=\sum_{j} c_{j} \mathfrak{g}_{j}$, where the coefficients $c_{j}$ are nonzero reals and the $\mathfrak{g}_{j}$ are log-free transmonomials. "Formal" means that we want to contemplate the sum as-is, not try to assign a "value" to it. The sum could even be transfinite (indexed by an ordinal), but for each term $c_{j} \mathfrak{g}_{j}$, the monomial $\mathfrak{g}_{j}$ is far smaller than all previous terms. Example:

$$
-4 e^{\sum_{j=0}^{\infty} x^{-j} e^{x}}+\sum_{j=0}^{\infty} x^{-j} e^{x}-17+\pi x^{-1}
$$

Transseries are added termwise (even series of transseries, but each monomial should occur only a finite number of times, so we can collect them). Transseries are multiplied in the way suggested by the notation"multiply it out" -but again we have to make sure that each monomial occurs in the product only a finite number of times. The transseries $T=\sum c_{j} \mathfrak{g}_{j}$ is purely large iff $\mathfrak{g}_{j} \succ 1$ for all terms $c_{j} \mathfrak{g}_{j}$; and $T$ is small iff $g_{j} \prec 1$ for all terms $c_{j} \mathfrak{g}_{j}$. A nonzero transseries $T=\sum c_{j} \mathfrak{g}_{j}$ has a dominant term $c_{0} \mathfrak{g}_{0}$ with $\mathfrak{g}_{0} \succ \mathfrak{g}_{j}$ for all other terms $c_{j} \mathfrak{g}_{j}$. If $c_{0}>0$ we say $T>0$. An ordering $>$ is then defined by: $S>T$ iff $S-T>0$.
We consider only transmonomials and transseries of "finite exponential height" -so, for example, these are not allowed:

$$
e^{e^{e^{e \cdot}+x}+x}+x, \quad x^{-1}+e^{-x}+e^{-e^{x}}+e^{-e^{e^{x}}}+\cdots
$$

(c) Differentiation is defined as in elementary calculus:

$$
\left(x^{b} e^{L}\right)^{\prime}=b x^{b-1} e^{L}+x^{b} L^{\prime} e^{L}, \quad\left(\sum c_{j} \mathfrak{g}_{j}\right)^{\prime}=\sum c_{j} \mathfrak{g}_{j}^{\prime}
$$

(d) Write $\log _{m} x$ for $\log \log \cdots \log x$ with $m \log$, where $m$ is a nonnegative integer. A general transseries is obtained by substitution of some $\log _{m} x$ for $x$ in a log-free transseries. Example:

$$
e^{(\log \log x)^{1 / 2}+x}+(\log \log x)^{1 / 2}+x^{-2}
$$

A general transmonomial is obtained similarly from a log-free transmonomial.

There are a few additional features in the development, as we will see in Section 3. But for now let's proceed to some examples. Computations with transseries can seem natural in many cases, even without the technical definitions. And-as with generating functions-even if they do not converge.

Example 2.1. Let us multiply $A=x-1$ times $B=\sum_{j=0}^{\infty} x^{-j}$.

$$
\begin{aligned}
(x & -1)\left(1+x^{-1}+x^{-2}+x^{-3}+\ldots\right) \\
& =x \cdot\left(1+x^{-1}+x^{-2}+x^{-3}+\ldots\right)-1 \cdot\left(1+x^{-1}+x^{-2}+x^{-3}+\ldots\right) \\
& =x+1+x^{-1}+x^{-2}+\cdots-1-x^{-1}-x^{-2}-x^{-3}-\cdots \\
& =x
\end{aligned}
$$

Example 2.2. Both transseries

$$
S=\sum_{j=0}^{\infty} j!x^{-j}, \quad T=\sum_{j=0}^{\infty}(-1)^{j} j!x^{-j}
$$

are divergent. For the product: the combinatorial identity

$$
\sum_{j=0}^{n}(-1)^{j} j!(n-j)!= \begin{cases}\frac{(n+1)!}{1+n / 2}, & n \text { even } \\ 0, & n \text { odd }\end{cases}
$$

means that

$$
S T=\sum_{j=0}^{\infty} \frac{(2 j+1)!}{j+1} x^{-2 j}
$$

Example 2.3. Now consider

$$
U=\sum_{j=1}^{\infty} j e^{-j x}, \quad V=\sum_{k=0}^{\infty} x^{-k}
$$

When $U V$ is multiplied out, each monomial $x^{-k} e^{-j x}$ occurs only once, so our result is a transseries whose support has order type $\omega^{2}$.

$$
U V=\sum_{j=1}^{\infty}\left(\sum_{k=0}^{\infty} j x^{-k} e^{-j x}\right)
$$

(For an explanation of order type, see [32, p. 27] or [57, p. 127] or even [59].) Example 2.4. Every nonzero transseries has a multiplicative inverse. What is the inverse of $e^{x}+x$ ? Use the Taylor series for $1 /(1+z)$ like this:

$$
\begin{aligned}
\left(e^{x}+x\right)^{-1} & =\left(e^{x}\left(1+x e^{-x}\right)\right)^{-1}=e^{-x} \sum_{j=0}^{\infty}(-1)^{j}\left(x e^{-x}\right)^{j} \\
& =\sum_{j=0}^{\infty}(-1)^{j} x^{j} e^{-(j+1) x}
\end{aligned}
$$

Example 2.5. The hyperbolic sine is a two-term transseries, $\sinh x=(1 / 2) e^{x}-$ $(1 / 2) e^{-x}$. Let's compute its logarithm. Use the Taylor series for $\log (1-z)$.

$$
\log (\sinh x)=\log \left(\frac{e^{x}}{2}\left(1-e^{-2 x}\right)\right)=x-\log 2-\sum_{j=1}^{\infty} \frac{e^{-2 j x}}{j}
$$

Wasn't that easy?
Example 2.6. How about the inverse of

$$
T=\sum_{j=0}^{\infty} j!x^{-j-1}=x^{-1}+x^{-2}+2 x^{-3}+6 x^{-4}+24 x^{-5}+\cdots ?
$$

We can compute as many terms as we want, with enough effort. First, $T=$ $x^{-1}(1+V)$, where $V=x^{-1}+2 x^{-2}+6 x^{-3}+24 x^{-4}+\cdots$ is small. So

$$
\begin{array}{rl}
T^{-1}= & \left(x^{-1}\right)^{-1}(1+V)^{-1}=x\left(1-V+V^{2}-V^{3}+V^{4}-\ldots\right) \\
=x & x\left[1-\left(x^{-1}+2 x^{-2}+6 x^{-3}+24 x^{-4}+\ldots\right)\right. \\
& \quad+\left(x^{-1}+2 x^{-2}+6 x^{-3}+\ldots\right)^{2} \\
& \left.\quad-\left(x^{-1}+2 x^{-2}+\ldots\right)^{3}+\left(x^{-1}+\ldots\right)^{4}+\ldots\right] \\
= & x-1-x^{-1}-3 x^{-2}-13 x^{-3}+\cdots .
\end{array}
$$

Searching the On-Line Encyclopedia of Integer Sequences [56] shows that these coefficients are sequence A003319.
Example 2.7. Function $x e^{x}$ has compositional inverse known as the Lambert $\boldsymbol{W}$ function. So $W(x) e^{W(x)}=x$. The transseries is:

$$
\begin{aligned}
W(x)= & \log x-\log \log x+\frac{\log \log x}{\log x}+\frac{(\log \log x)^{2}}{2(\log x)^{2}}-\frac{\log \log x}{(\log x)^{2}}+\frac{(\log \log x)^{3}}{3(\log x)^{3}} \\
& -\frac{3(\log \log x)^{2}}{2(\log x)^{3}}+\frac{\log \log x}{(\log x)^{3}}+\frac{(\log \log x)^{4}}{4(\log x)^{4}}-\frac{11(\log \log x)^{3}}{6(\log x)^{4}}+\cdots
\end{aligned}
$$

We will see below (Problem 4.2) how to compute this. But for now, let's see how to compute $e^{W(x)}$. The two terms $\log x$ and $\log \log x$ are large, the rest is small. If $W(x)=\log x-\log \log x+S$, then

$$
e^{W(x)}=e^{\log x} e^{-\log \log x} e^{S}=\frac{x}{\log x}\left(\sum_{j=0}^{\infty} \frac{S^{j}}{j!}\right)
$$

Then put in $S=\log \log x / \log x+\cdots$, as many terms as needed, to get

$$
e^{W(x)}=\frac{x}{\log x}+\frac{x \log \log x}{(\log x)^{2}}+\frac{x(\log \log x)^{2}}{(\log x)^{3}}-\frac{x \log \log x}{(\log x)^{3}}+\cdots
$$

This is $e^{W(x)}$. Now we can multiply this by the original $W$ :

$$
W(x) e^{W(x)}=x+\cdots
$$

where the missing terms are of order higher than computed. In fact, the claim is that all higher terms cancel.
Remark 2.8. By a general result of van den Dries-Macintyre-Marker (3.12 and 6.30 in [18]), there exists a coherent way to associate a transseries expansion at $+\infty$ to every function $(a,+\infty) \rightarrow \mathbb{R}$ (where $a \in \mathbb{R}$ ) which, like the functions considered in Examples 2.4 to 2.7, is definable (in the sense of mathematical logic) from real constants, addition, multiplication, and exp.

## Écalle-Borel Summation.

There is a system to assign real functions to many transseries. Écalle calls this "accelero-summation". It is a vast generalization of the classical Borel summation method. Here we will consider transseries only as formal objects, for the most part, but I could not resist including a few remarks on summation.

The basic Borel summation works like this: The Laplace transform $\mathcal{L}$ is defined by

$$
\mathcal{L}[F](x)=\int_{0}^{\infty} e^{-x p} F(p) d p
$$

when it exists. The inverse Laplace transform, or Borel transform, will be written $\mathcal{B}$, so that $\mathcal{B}[f]=F$ iff $\mathcal{L}[F]=f$. (E. Post [48] provided an interesting formula for $\mathcal{B}$; see also [39], [7].) The composition $\mathcal{L B}$ is an "isomorphism" in the sense that it preserves "all operations"-whatever that means; perhaps in the wishful sense. In fact, in some cases even if $f$ is merely a formal series (a divergent series), still $\mathcal{L B}[f]$ yields an actual function. If so, that is the Borel sum of the series.

We will use variable $x$ in physical space, and variable $p$ in Borel space. Then compute $\mathcal{L}\left[p^{n}\right]=n!x^{-n-1}$ for $n \in \mathbb{N}$, so $\mathcal{B}\left[x^{-j}\right]=p^{j-1} /(j-1)$ ! for integers $j \geq 1$.
Example 2.9. Borel summation works on the series $f=\sum_{j=0}^{\infty} 2^{j} x^{-j}$. (Except for the first term-no delta functions here.) Write $f=1+g$. First $\mathcal{B}[g]=$ $\sum_{j=1}^{\infty} 2^{j} p^{j-1} /(j-1)!=2 e^{2 p}$. Then

$$
\mathcal{L B}[g](x)=\int_{0}^{\infty} 2 e^{-x p} e^{2 p} d p=\frac{2}{x-2}
$$

Adding the 1 back on, we conclude that the sum of the series should be

$$
1+\frac{2}{x-2}=\frac{x}{x-2}=\frac{1}{1-2 / x}
$$

as expected.
Of course the formal series $f=\sum_{j=0}^{\infty} 2^{j} x^{-j}$ satisfies $f \cdot(1-2 / x)=1$. So if $\mathcal{L B}$ is supposed to preserve all operations then there is no other sum possible.
Example 2.10. Consider Euler's series $f=\sum_{j=0}^{\infty}(-1)^{j} j!x^{-j-1}$, a series that diverges for all $x$. So we want: $\mathcal{B}[f]=\sum_{j=0}^{\infty}(-1)^{j} p^{j}=1 /(1+p)$. This expression makes sense for all $p \geq 0$, not just the ones within the radius of convergence. Then $\mathcal{B}[f]$ should be $1 /(1+p)$. Then the Laplace integral converges,

$$
\mathcal{L B}[f](x)=\int_{0}^{\infty} \frac{e^{-x p} d p}{1+p}=-e^{x} \operatorname{Ei}(-x)
$$

This is the Borel sum of the series $f$.
Similarly, consider the series

$$
g=\sum_{j=0}^{\infty} \frac{j!}{x^{j+1}} .
$$

In the same way, we get

$$
\mathcal{L B}[g](x)=\int_{0}^{\infty} \frac{e^{-x p} d p}{1-p}
$$

where now (because of the pole at $p=1$ ) this taken as a principal value integral, and we get $e^{-x} \operatorname{Ei}(x)$ as the value.

Borel summation is the beginning of the story. Much more powerful methods have been developed. (Écalle invented most of the techniques, then others have made them rigorous and improved them.) To a large extent it is known that transseries that arise (from ODEs, PDEs, difference equations, etc.) can be summed, and much more is suspected. This summation is virtually as faithful as convergent summation. But the subject is beyond the scope of this paper. In fact, it seems that a simple exposition is not possible with our present understanding. For more on summation see [11, §3.1], [12], [13].

## 3 The Formal Construction.

Now we come to the technical part of the paper. All a mathematician needs to know in order to apply transseries.

To do the types of computations we have seen, a formal construction is desirable. It should allow not only "formal power series," but also exponentials and logarithms. In reading this, you can note that in fact we are not really using high-level mathematics.

Descriptions of the system of transseries are found, for example, in $[1,13$, $18,37]$. But those accounts are (to a greater or lesser extent) technical and involve jargon of the subfield. It is hoped that by carefully reading this section, a reader who is not a specialist will be able to understand the simplicity of the construction. Some details are not checked here, especially the tedious ones.

Items called Comment, enclosed between two $\circ$ signs, are not part of the formal construction. They are included as illustration and motivation. Perhaps these commentaries cannot be completely understood until after the formal construction has been read.

- Comment 3.1. Functions (or expressions) of the form $x^{a} e^{b x}$, where $a, b \in \mathbb{R}$, are transmonomials. (There are also many other transmonomials. But these will be enough for most of our illustrative comments.) We may think of the "far larger" relation $\succ$ describing relative size when $x \rightarrow+\infty$. In particular, $x^{a_{1}} e^{b_{1} x} \succ x^{a_{2}} e^{b_{2} x}$ if and only if $b_{1}>b_{2}$ or $\left\{b_{1}=b_{2}\right.$ and $\left.a_{1}>a_{2}\right\}$.


## 3A Multi-Indices.

- Comment 3.2. The set $\mathfrak{G}$ of monomials is a group under multiplication. This group (even the subgroup of monomials $x^{a} e^{b x}$ ) is not finitely generated. But sometimes we will want to consider a finitely generated subgroup of $\mathfrak{G}$. If $\mu_{1}, \ldots, \mu_{n}$ is a set of generators, then the generated group is

$$
\left\{\mu_{1}^{k_{1}} \mu_{2}^{k_{2}} \cdots \mu_{n}^{k_{n}}: k_{1}, k_{2}, \ldots, k_{n} \in \mathbb{Z}\right\}
$$

We will discuss the use of multi-indices $\mathbf{k}=\left(k_{1}, k_{2}, \ldots, k_{n}\right)$ so that later $\mu_{1}^{k_{1}} \mu_{2}^{k_{2}} \cdots \mu_{n}^{k_{n}}$ can be abbreviated $\boldsymbol{\mu}^{\mathbf{k}}$ and save much writing.

It does no harm to omit the group identity 1 from a list of generators; replacing some generators by their inverses, we may assume the generators $\mu_{j}$ are all small: $\mu_{j} \prec 1$. (We will think of these as "ratios" between one term of a series and the next. A ratio set is a finite set of small monomials.) Then the correspondence between multi-indices $\mathbf{k}$ and monomials $\boldsymbol{\mu}^{\mathbf{k}}$ reverses the ordering. (That is, if $\mathbf{k}>\mathbf{p}$, then $\boldsymbol{\mu}^{\mathbf{k}} \prec \boldsymbol{\mu}^{\mathbf{p}}$.) This means terminology that seems right on one side may seem to be backward on the other side. Even with conventional asymptotic series, larger terms are written to the left, smaller terms to the right, reversing the convention for a number line.

Begin with a positive integer $n$. The set $\mathbb{Z}^{n}$ of $n$-tuples of integers is a group under componentwise addition. For notation-avoiding subscripts, since we want to use subscripts for many other things-if $\mathbf{k} \in \mathbb{Z}^{n}$ and $1 \leq i \leq n$, write $\mathbf{k}[i]$ for the $i$ th component of $\mathbf{k}$. The partial order $\leq$ is defined by: $\mathbf{k} \leq \mathbf{p}$ iff $\mathbf{k}[i] \leq \mathbf{p}[i]$ for all $i$. And $\mathbf{k}<\mathbf{p}$ iff $\mathbf{k} \leq \mathbf{p}$ and $\mathbf{k} \neq \mathbf{p}$. Element $\mathbf{0}=(0,0, \ldots, 0)$ is the identity for addition.
Definition 3.3. For $\mathbf{m} \in \mathbb{Z}^{n}$, define $\mathbf{J}_{\mathbf{m}}=\left\{\mathbf{k} \in \mathbb{Z}^{n}: \mathbf{k} \geq \mathbf{m}\right\}$.

- Comment 3.4. For example $\mathbf{J}_{(-1,2)}=\left\{(k, l) \in \mathbb{Z}^{2}: k \geq-1, l \geq 2\right\}$. The sets $\mathbf{J}_{\mathbf{m}}$ will be used below (Definition 3.34) to define "grids" of monomials. If $\mu_{1}=x^{-1}$ and $\mu_{2}=e^{-x}$ are the ratios making up the ratio set $\boldsymbol{\mu}$, then

$$
\left\{\boldsymbol{\mu}^{\mathbf{k}}: \mathbf{k} \in \mathbf{J}_{(-1,2)}\right\}=\left\{x^{-k} e^{-l x}: k \geq-1, l \geq 2\right\}
$$

is the corresponding grid.
Write $\mathbb{N}=\{0,1,2,3, \cdots\}$ including 0 . The subset $\mathbb{N}^{n}$ of $\mathbb{Z}^{n}$ is closed under addition. Note $\mathbf{J}_{\mathbf{m}}$ is the translate of $\mathbb{N}^{n}$ by $\mathbf{m}$. That is, $\mathbf{J}_{\mathbf{m}}=$ $\left\{\mathbf{k}+\mathbf{m}: \mathbf{k} \in \mathbb{N}^{n}\right\}$. And $\mathbb{N}^{n}=\mathbf{J}_{\mathbf{0}}$. Translation preserves order.

The next three propositions explain that the set $\mathbf{J}_{\mathbf{m}}$ is well-partially-ordered (also called Noetherian). These three - which are collectively known as "Dickson's Lemma" - are the main reason why certain algorithms in computer algebra (Gröbner bases) terminate. Equivalent properties and the usefulness in power series rings are discussed in Higman [33] and Erdős \& Radó (cited in [33]).
Proposition 3.5. If $\mathbf{E} \subseteq \mathbf{J}_{\mathbf{m}}$ and $\mathbf{E} \neq \varnothing$, then there is a minimal element: $\mathbf{k}_{0} \in \mathbf{E}$ and $\mathbf{k}<\mathbf{k}_{0}$ holds for no element $\mathbf{k} \in \mathbf{E}$.
Proof. Because translation preserves order, it suffices to do the case of $\mathbf{J}_{\mathbf{0}}=$ $\mathbb{N}^{n}$. First, $\{\mathbf{k}[1]: \mathbf{k} \in \mathbf{E}\}$ is a nonempty subset of $\mathbb{N}$, so it has a least element, say $m_{1}$. Then $\left\{\mathbf{k}[2]: \mathbf{k} \in \mathbf{E}, \mathbf{k}[1]=m_{1}\right\}$ is a nonempty subset of $\mathbb{N}$, so it has a least element, say $m_{2}$. Continue. Then $\mathbf{k}_{0}=\left(m_{1}, \ldots, m_{n}\right)$ is minimal in E.

Proposition 3.6. Let $\mathbf{E} \subseteq \mathbf{J}_{\mathbf{m}}$ be infinite. Then there is a sequence $\mathbf{k}_{j} \in \mathbf{E}$, $j \in \mathbb{N}$, with $\mathbf{k}_{0}<\mathbf{k}_{1}<\mathbf{k}_{2}<\cdots$.

Proof. It is enough to do the case $\mathbb{N}^{n}$. The proof is by induction on $n$-it is true for $n=1$. Assume $n \geq 2$. Define the set $\widetilde{\mathbf{E}} \subseteq \mathbb{Z}^{n-1}$ by

$$
\widetilde{\mathbf{E}}=\{(\mathbf{k}[1], \mathbf{k}[2], \ldots, \mathbf{k}[n-1]): \mathbf{k} \in \mathbf{E}\} .
$$

Case 1. $\widetilde{\mathbf{E}}$ is finite. Then for some $\mathbf{p} \in \widetilde{\mathbf{E}}$, the set

$$
\mathbf{E}^{\prime}=\{k \in \mathbb{N}:(\mathbf{p}[1], \ldots, \mathbf{p}[n-1], k) \in \mathbf{E}\}
$$

is infinite. Choose an increasing sequence $k_{j} \in \mathbf{E}^{\prime}$ to get the increasing sequence in $\mathbf{E}$.

Case 2. $\widetilde{\mathbf{E}}$ is infinite. By the induction hypothesis, there is a strictly increasing sequence $\mathbf{p}_{j} \in \widetilde{\mathbf{E}}$. So there is a sequence $\mathbf{k}_{j} \in \mathbf{E}$ that is increasing in every coordinate except possibly the last. If some last coordinate occurs infinitely often, use it to get an increasing sequence in $E$. If not, choose a subsequence of these last coordinates that increases.

Proposition 3.7. Let $\mathbf{E} \subseteq \mathbf{J}_{\mathbf{m}}$. Then the set Min $\mathbf{E}$ of all minimal elements of $\mathbf{E}$ is finite. For every $\mathbf{k} \in \mathbf{E}$, there is $\mathbf{k}_{0} \in \operatorname{Min} \mathbf{E}$ with $\mathbf{k}_{0} \leq \mathbf{k}$.
Proof. No two minimal elements are comparable, so $\operatorname{Min} \mathbf{E}$ is finite by Proposition 3.6. If $\mathbf{E}=\varnothing$, then $\operatorname{Min} \mathbf{E}=\varnothing$ vacuously satisfies this. Suppose $\mathbf{E} \neq \varnothing$. Then $\operatorname{Min} \mathbf{E} \neq \varnothing$ satisfies the required conclusion by Proposition 3.5.

## Convergence of sets

Write $\triangle$ for the symmetric difference operation on sets. We will define convergence of a sequence of sets $\mathbf{E}_{j} \subseteq \mathbb{Z}^{n}$ (or indeed any infinite collection $\left(\mathbf{E}_{i}\right)_{i \in I}$ of sets).
Definition 3.8. Let $I$ be an infinite index set, and for each $i \in I$, let $\mathbf{E}_{i} \subseteq \mathbb{Z}^{n}$ be given. We say the family $\left(\mathbf{E}_{i}\right)_{i \in I}$ is point-finite iff each $\mathbf{p} \in \mathbb{Z}^{n}$ belongs to $\mathbf{E}_{i}$ for only finitely many $i$. Let $\mathbf{m} \in \mathbb{Z}^{n}$. We write $\mathbf{E}_{i} \xrightarrow{\mathbf{m}} \varnothing$ iff $\mathbf{E}_{i} \subseteq \mathbf{J}_{\mathbf{m}}$ for all $i$ and $\left(\mathbf{E}_{i}\right)$ is point-finite. We write $\mathbf{E}_{i} \rightarrow \varnothing$ iff there exists $\mathbf{m}$ such that $\mathbf{E}_{i} \xrightarrow{\mathbf{m}} \varnothing$. Furthermore, write $\mathbf{E}_{i} \xrightarrow{\mathbf{m}} \mathbf{E}$ iff $\mathbf{E}_{i} \subseteq \mathbf{J}_{\mathbf{m}}$ for all $i$ and $\mathbf{E}_{i} \triangle \mathbf{E} \xrightarrow{\mathbf{m}} \varnothing$; and write $\mathbf{E}_{i} \rightarrow \mathbf{E}$ iff $\mathbf{E}_{i} \xrightarrow{\mathbf{m}} \mathbf{E}$ for some $\mathbf{m}$.

Expressed differently: $\left(\mathbf{E}_{i} \triangle \mathbf{E}\right)_{i \in I}$ is point-finite if and only if $\lim \sup _{i} \mathbf{E}_{i}=$ $\liminf \inf _{i}=\mathbf{E}$.

- Comment 3.9. Examples in $\mathbb{Z}=\mathbb{Z}^{1}$. Let $\mathbf{E}_{i}=\{i, i+1, i+2, \cdots\}$ for $i \in \mathbb{N}$. Then the sequence $\mathbf{E}_{i}$ is point-finite. And $\mathbf{E}_{i} \rightarrow \varnothing$. But let $\mathbf{F}_{i}=\{-i\}$ for $i \in \mathbb{N}$. Again the sequence $\mathbf{F}_{i}$ is point-finite, but there is no $\mathbf{m}$ with $\mathbf{F}_{i} \subseteq \mathbf{J}_{\mathbf{m}}$ for all $i$, so $\mathbf{F}_{i}$ does not converge in this sense.

This type of convergence is metrizable when restricted to any $\mathbf{J}_{\mathbf{m}}$.
Notation 3.10. For $\mathbf{k}=\left(k_{1}, k_{2}, \ldots, k_{n}\right)$, define $|\mathbf{k}|=k_{1}+k_{2}+\cdots+k_{n}$.
Proposition 3.11. Let $\mathbf{m} \in \mathbb{Z}^{n}$. For $\mathbf{E}, \mathbf{F} \subseteq \mathbf{J}_{\mathbf{m}}$, define

$$
d(\mathbf{E}, \mathbf{F})=\sum_{\mathbf{k} \in \mathbf{E} \triangle \mathbf{F}} 2^{-|\mathbf{k}|} .
$$

Then for any sets $\mathbf{E}_{i} \subseteq \mathbf{J}_{\mathbf{m}}$, we have $\mathbf{E}_{i} \rightarrow \mathbf{E}$ if and only if $d\left(\mathbf{E}_{i}, \mathbf{E}\right) \rightarrow 0$. And $d$ is a metric on subsets of $\mathbf{J}_{\mathbf{m}}$.

## 3B Hahn Series.

We begin with an ordered abelian group $\mathfrak{M}$, called the monomial group (or valuation group). By "ordered" we mean totally ordered or linearly ordered. The operation is written multiplicatively, the identity is 1 , the order relation is $\succ$ and read "far larger than". This is a "strict" order relation; that is, $\mathfrak{g} \succ \mathfrak{g}$ is false. An element $\mathfrak{g} \in \mathfrak{M}$ is called large iff $\mathfrak{g} \succ 1$, and small if $\mathfrak{g} \prec 1$. [We will use Fraktur letters: lower case for monomials and upper case for sets of monomials.]

- Comment 3.12. The material in Subsections 3B and 3C will apply to any ordered abelian group $\mathfrak{M}$. Later (Subsections 3D and 3E) we will construct the particular group that will specialize this general construction into the transseries construction. Comments will use the group of monomials $x^{a} e^{b x}$ discussed above.

We use the field $\mathbb{R}$ of real numbers for values. Write $\mathbb{R}^{\mathfrak{M}}$ for the set of functions $T: \mathfrak{M} \rightarrow \mathbb{R}$. For $T \in \mathbb{R}^{\mathfrak{M}}$ and $\mathfrak{g} \in \mathfrak{M}$, we will use square brackets $T[\mathfrak{g}]$ for the value of $T$ at $\mathfrak{g}$, because later we will want to use round brackets $T(x)$ in another more common sense.
Definition 3.13. The support of a function $T \in \mathbb{R}^{\mathfrak{M}}$ is

$$
\operatorname{supp} T=\{\mathfrak{g} \in \mathfrak{M}: T[\mathfrak{g}] \neq 0\}
$$

Let $\mathfrak{A} \subseteq \mathfrak{M}$. We say $T$ is supported by $\mathfrak{A}$ if $\operatorname{supp} T \subseteq \mathfrak{A}$.
Notation 3.14. In fact, $T$ will usually be written as a formal combination of group elements. That is:

$$
T=\sum_{\mathfrak{g} \in \mathfrak{A}} a_{\mathfrak{g}} \mathfrak{g}, \quad a_{\mathfrak{g}} \in \mathbb{R}
$$

will be used for the function $T$ with $T[\mathfrak{g}]=a_{\mathfrak{g}}$ for $\mathfrak{g} \in \mathfrak{A}$ and $T[\mathfrak{g}]=0$ otherwise. The set $\mathfrak{A}$ might or might not be the actual support of $T$. Accordingly, such $T$ may be called a Hahn series or generalized power series.
Definition 3.15. If $c \in \mathbb{R}$, then $c 1 \in \mathbb{R}^{\mathfrak{M}}$ is called a constant and identified with $c$. (That is, $T[1]=c$ and $T[\mathfrak{g}]=0$ for all $\mathfrak{g} \neq 1$.) If $\mathfrak{m} \in \mathfrak{M}$, then $1 \mathfrak{m} \in \mathbb{R}^{\mathfrak{M}}$ is called a monomial and identified with $\mathfrak{m}$. (That is, $T[\mathfrak{m}]=1$ and $T[\mathfrak{g}]=0$ for all $\mathfrak{g} \neq \mathfrak{m}$.)

In all cases of interest to us, the support will be well ordered (according to the converse of $\succ$ ). That is, for all $\mathfrak{A} \subseteq \operatorname{supp}(T)$, if $\mathfrak{A} \neq \varnothing$, it has a maximum: $\mathfrak{m} \in \mathfrak{A}$ such that for all $\mathfrak{g} \in \mathfrak{A}$, if $\mathfrak{g} \neq \mathfrak{m}$, then $\mathfrak{m} \succ \mathfrak{g}$. We will use the term "well ordered" and not "converse well ordered" for this, so the reader should be wary.

Proposition 3.16. Let $\mathfrak{A} \subseteq \mathfrak{M}$ be well ordered for the converse of $\succ$. Every infinite subset in $\mathfrak{A}$ contains an infinite strictly decreasing sequence $\mathfrak{g}_{1} \succ \mathfrak{g}_{2} \succ$ $\cdots$. There is no infinite strictly increasing sequence in $\mathfrak{A}$.

Definition 3.17. Let $T \neq 0$ be

$$
T=\sum_{\mathfrak{g} \in \mathfrak{A}} a_{\mathfrak{g}} \mathfrak{g}, \quad a_{\mathfrak{g}} \in \mathbb{R}
$$

with $\mathfrak{m} \in \mathfrak{A}, \mathfrak{m} \succ \mathfrak{g}$ for all other $\mathfrak{g} \in \mathfrak{A}$, and $a_{\mathfrak{m}} \neq 0$. Then the magnitude of $T$ is $\operatorname{mag} T=\mathfrak{m}$, the leading coefficient of $T$ is $a_{\mathfrak{m}}$, and the dominance of $T$ is $\operatorname{dom} T=a_{\mathfrak{m}} \mathfrak{m}$. We say $T$ is positive if $a_{\mathfrak{m}}>0$ and write $T>0$. We say $T$ is negative if $a_{\mathfrak{m}}<0$ and write $T<0$. We say $T$ is small if $\mathfrak{g} \prec 1$ for all $\mathfrak{g} \in \operatorname{supp} T$ (equivalently: $\operatorname{mag} T \prec 1$ or $T=0$ ). We say $T$ is large if $\operatorname{mag} T \succ 1$. We say $T$ is purely large if $\mathfrak{g} \succ 1$ for all $\mathfrak{g} \in \operatorname{supp} T$. (Because of the standard empty-set conventions: 0, although not large, is purely large.)

Remark 3.18. Alternate terminology [37, 54]: magnitude = leading monomial; dominance $=$ leading term; large $=$ infinite; small $=$ infinitesimal.

- Comment 3.19. Let $A=-3 e^{x}+4 x^{2}$. Then mag $A=e^{x}$, $\operatorname{dom} A=-3 e^{x}, A$ is negative, $A$ is large, $A$ is purely large.

Definition 3.20. Addition is defined by components: $(A+B)[\mathfrak{g}]=A[\mathfrak{g}]+B[\mathfrak{g}]$. Constant multiples $c A$ are also defined by components.

Remark 3.21. The union of two well ordered sets is well ordered. So if $A, B$ each have well ordered support, so does $A+B$.
Notation 3.22. We say $S>T$ if $S-T>0$. For nonzero $S$ and $T$ we say $S \succ T$ (read $S$ is far larger than $T)$ iff $\operatorname{mag} S \succ \operatorname{mag} T$; we say $S \asymp T(\operatorname{read} S$ is comparable to $T$ or $S$ has the same magnitude as $T$ ) iff mag $S=\operatorname{mag} T$; and we say $S \sim T(\operatorname{read} S$ is asymptotic to $T)$ iff $\operatorname{dom} S=\operatorname{dom} T$. Write $S \succcurlyeq T$ iff $S \succ T$ or $S \asymp T$. These have reversed counterparts as usual: for example, $S<T$ means $T>S$.

- Comment 3.23. Examples:

$$
\begin{aligned}
& -3 e^{x}+4 x^{2}<x^{9} \\
& -3 e^{x}+4 x^{2} \succ x^{9} \\
& -3 e^{x}+4 x^{2} \asymp 7 e^{x}+x^{9} \\
& -3 e^{x}+4 x^{2} \sim-3 e^{x}+x^{9}
\end{aligned}
$$

- 


## The Two Canonical Decompositions

Proposition 3.24 (Canonical Additive Decomposition). Every A may be written uniquely in the form $A=L+c+S$, where $L$ is purely large, $c$ is a constant, and $S$ is small.

Remark 3.25. Terminology: $L$ is the purely large part, $c$ is the constant term, and $S$ is the small part of $A$.
Definition 3.26. Multiplication is defined by convolution (as suggested by the formal sum notation).

$$
\begin{aligned}
& \sum_{\mathfrak{g} \in \mathfrak{M}} a_{\mathfrak{g}} \mathfrak{g} \cdot \sum_{\mathfrak{g} \in \mathfrak{M}} b_{\mathfrak{g}} \mathfrak{g}=\sum_{\mathfrak{g} \in \mathfrak{M}}\left(\sum_{\mathfrak{m} \mathfrak{n}=\mathfrak{g}} a_{\mathfrak{m}} b_{\mathfrak{n}}\right) \mathfrak{g} \\
& \text { or } \quad(A B)[\mathfrak{g}]=\sum_{\mathfrak{m} \mathfrak{n}=\mathfrak{g}} A[\mathfrak{m}] B[\mathfrak{n}]
\end{aligned}
$$

Products are defined at least for $A, B$ with well ordered support.
Proposition 3.27. If $\mathfrak{A}_{1}, \mathfrak{A}_{2} \subseteq \mathfrak{M}$ are well ordered sets (for the converse of $\succ)$, then $\mathfrak{A}=\left\{\mathfrak{g}_{1} \mathfrak{g}_{2}: \mathfrak{g}_{1} \in \mathfrak{A}_{1}, \mathfrak{g}_{2} \in \mathfrak{A}_{2}\right\}$ is also well ordered. For every $\mathfrak{g} \in \mathfrak{A}$, the set

$$
\left\{\left(\mathfrak{g}_{1}, \mathfrak{g}_{2}\right): \mathfrak{g}_{1} \in \mathfrak{A}_{1}, \mathfrak{g}_{2} \in \mathfrak{A}_{2}, \mathfrak{g}_{1} \mathfrak{g}_{2}=\mathfrak{g}\right\}
$$

is finite.
Proof. Let $\mathfrak{B} \subseteq \mathfrak{A}$ be nonempty. Assume $\mathfrak{B}$ has no greatest element. Then there exist sequences $\mathfrak{m}_{j} \in \mathfrak{A}_{1}$ and $\mathfrak{n}_{j} \in \mathfrak{A}_{2}$ with $\mathfrak{m}_{j} \mathfrak{n}_{j} \in \mathfrak{B}$ and $\mathfrak{m}_{1} \mathfrak{n}_{1} \prec$ $\mathfrak{m}_{2} \mathfrak{n}_{2} \prec \cdots$. Because $\mathfrak{A}_{1}$ is well ordered, taking a subsequence we may assume $\mathfrak{m}_{1} \succcurlyeq \mathfrak{m}_{2} \succcurlyeq \cdots$. But then $\mathfrak{n}_{1} \prec \mathfrak{n}_{2} \prec \cdots$, so $\mathfrak{A}_{2}$ is not well ordered.

Suppose $\left(\mathfrak{g}_{1}, \mathfrak{g}_{2}\right),\left(\mathfrak{m}_{1}, \mathfrak{m}_{2}\right) \in \mathfrak{A}_{1} \times \mathfrak{A}_{2}$ with $\mathfrak{g}_{1} \mathfrak{g}_{2}=\mathfrak{g}=\mathfrak{m}_{1} \mathfrak{m}_{2}$. If $\mathfrak{g}_{1} \neq \mathfrak{m}_{1}$, then $\mathfrak{g}_{2} \neq \mathfrak{m}_{2}$. If $\mathfrak{g}_{1} \succ \mathfrak{m}_{1}$, then $\mathfrak{g}_{2} \prec \mathfrak{m}_{2}$. Any infinite subset of a well ordered set contains an infinite strictly decreasing sequence, but the other well ordered set contains no infinite strictly increasing sequence.

Proposition 3.28. The set of all $T \in \mathbb{R}^{\mathfrak{M}}$ with well ordered support is an (associative, commutative) algebra over $\mathbb{R}$ with the operations defined above.

There are a lot of details to check. In fact this is a field [9, p. 276], but we won't need that result. This goes back to H. Hahn, 1907 [29].

Proposition 3.29 (Canonical Multiplicative Decomposition). Every nonzero $T \in \mathbb{R}^{\mathfrak{M}}$ with well ordered support may be written uniquely in the form $T=$ $a \cdot \mathfrak{g} \cdot(1+S)$ where $a$ is nonzero real, $\mathfrak{g} \in \mathfrak{M}$, and $S$ is small.

- Comment 3.30. $-3 e^{x}+4 x^{2}=-3 \cdot e^{x} \cdot\left(1-(4 / 3) x^{2} e^{-x}\right)$.
- 

Proposition 3.31. The set of all purely large $T$ (including 0) is a group under addition. The set of all small $T$ is a group under addition. The set of all purely large $T$ (with well ordered support) is closed under multiplication. The set of all small $T$ (with well ordered support) is closed under multiplication.

## 3C Grids.

Some definitions will depend on a finite set of ratios (or generators). We will keep track of the set of ratios more than is customary. But it is useful for the proofs, and especially for the Fixed-Point Theorem 4.22.

Write $\mathfrak{M}^{\text {small }}=\{\mathfrak{g} \in \mathfrak{M}: \mathfrak{g} \prec 1\}$. A ratio set (or generating set) is a finite subset $\boldsymbol{\mu} \subset \mathfrak{M}^{\text {small }}$. We will use bold Greek for ratio sets. If convenient, we may number the elements of $\boldsymbol{\mu}$ in order, $\mu_{1} \succ \mu_{2} \succ \cdots \succ \mu_{n}$ and then consider $\boldsymbol{\mu}$ an ordered $n$-tuple.
Notation 3.32. Let $\boldsymbol{\mu}=\left\{\mu_{1}, \ldots, \mu_{n}\right\} \subset \mathfrak{M}^{\text {small }}$. For any multi-index $\mathbf{k}=$ $\left(k_{1}, \ldots, k_{n}\right) \in \mathbb{Z}^{n}$, define $\boldsymbol{\mu}^{\mathbf{k}}=\mu_{1}^{k_{1}} \cdots \mu_{n}^{k_{n}}$.

If $\mathbf{k}>\mathbf{p}$, then $\boldsymbol{\mu}^{\mathbf{k}} \prec \boldsymbol{\mu}^{\mathbf{p}}$. Also $\boldsymbol{\mu}^{\mathbf{0}}=1$. If $\mathbf{k}>\mathbf{0}$ then $\boldsymbol{\mu}^{\mathbf{k}} \prec 1$ (but not in general conversely).

- Comment 3.33. Let $\boldsymbol{\mu}=\left\{x^{-1}, e^{-x}\right\}$. Then $1 \succ \mu_{1}^{-1} \mu_{2}=x e^{-x}$, even though $(-1,1) \ngtr(0,0)$.
Definition 3.34. For ratio set $\boldsymbol{\mu}=\left\{\mu_{1}, \ldots, \mu_{n}\right\}$, let $\mathfrak{J}^{\mu}=\left\{\boldsymbol{\mu}^{\mathbf{k}}: \mathbf{k} \in \mathbb{Z}^{n}\right\}$, the subgroup generated by $\boldsymbol{\mu}$. If $\mathbf{m} \in \mathbb{Z}^{n}$, then define a subset of $\mathfrak{J}^{\mu}$ by

$$
\mathfrak{J}^{\mu, \mathbf{m}}=\left\{\boldsymbol{\mu}^{\mathbf{k}}: \mathbf{k} \in \mathbf{J}_{\mathbf{m}}\right\}=\left\{\boldsymbol{\mu}^{\mathbf{k}}: \mathbf{k} \geq \mathbf{m}\right\}
$$

The sets $\mathfrak{J}^{\boldsymbol{\mu}, \mathbf{m}}$ are called grids. A Hahn series $T \in \mathbb{R}^{\mathfrak{M}}$ supported by some grid is said to be grid-based. A set $\mathfrak{A} \subseteq \mathfrak{M}$ is called a subgrid if $\mathfrak{A} \subseteq \mathfrak{J}^{\boldsymbol{\mu}, \mathbf{m}}$ for some $\boldsymbol{\mu}, \mathbf{m}$.

Proposition 3.35. Let $\mathbf{W}$ be the set of all subgrids.
(a) $\mathfrak{J}^{\boldsymbol{\mu}, \mathbf{m}}$ is well ordered (by the converse of $\succ$ ).
(b) If $\boldsymbol{\mu} \subseteq \tilde{\boldsymbol{\mu}}$, then $\mathfrak{J}^{\boldsymbol{\mu}, \mathbf{m}} \subseteq \tilde{\mathfrak{J}}^{\tilde{\boldsymbol{\mu}}, \tilde{\mathbf{m}}}$ for some $\widetilde{\mathbf{m}}$.
(c) If $\mathfrak{A}, \mathfrak{B} \in \mathbf{W}$, then $\mathfrak{A} \cup \mathfrak{B} \in \mathbf{W}$.
(d) If $\mathfrak{A}, \mathfrak{B} \in \mathbf{W}$, then $\mathfrak{A} \cdot \mathfrak{B} \in \mathbf{W}$, where $\mathfrak{A} \cdot \mathfrak{B}:=\{\mathfrak{a} \mathfrak{b}: \mathfrak{a} \in \mathfrak{A}, \mathfrak{b} \in \mathfrak{B}\}$.

Proof. (a) Let $\mathfrak{B} \subseteq \mathfrak{J}^{\boldsymbol{\mu}, \mathbf{m}}$ be nonempty. Define $\mathbf{E}=\left\{\mathbf{k} \in \mathbf{J}_{\mathbf{m}}: \boldsymbol{\mu}^{\mathbf{k}} \in \mathfrak{B}\right\}$. Then the set $\operatorname{Min} \mathbf{E}$ of minimal elements of $\mathbf{E}$ is finite. So the greatest element of $\mathfrak{B}$ is $\max \left\{\boldsymbol{\mu}^{\mathbf{k}}: \mathbf{k} \in \operatorname{Min} \mathbf{E}\right\}$.
(b) Insert 0s for the extra entries of $\widetilde{\mathbf{m}}$.
(c) Use the union of the two $\mu \mathrm{s}$ and the minimum of the two ms.
(d) Use the union of the two $\boldsymbol{\mu}$ s and the sum of the two $\mathbf{m s}$.

Remark 3.36. By (c) and (d), if $S, T \in \mathbb{R}^{\mathfrak{M}}$ each have support in $\mathbf{W}$, then $S+T$ and $S T$ also have support in $\mathbf{W}$.
Remark 3.37. Write $\boldsymbol{\mu}=\left\{\mu_{1}, \mu_{2}, \ldots, \mu_{n}\right\}$ and $\mathbf{m}=\left(m_{1}, m_{2}, \ldots, m_{n}\right)$. Saying $T$ is supported by the grid $\mathfrak{J}^{\boldsymbol{\mu}, \mathbf{m}}$ means that $T=\sum c_{\mathbf{k}} \boldsymbol{\mu}^{\mathbf{k}}$ is a one-sided multiple Laurent series in the symbols $\boldsymbol{\mu}$ :

$$
\sum_{k_{1}=m_{1}}^{\infty} \sum_{k_{2}=m_{2}}^{\infty} \ldots \sum_{k_{n}=m_{n}}^{\infty} c_{k_{1} k_{2} \ldots k_{n}} \mu_{1}^{k_{1}} \mu_{2}^{k_{2}} \cdots \mu_{n}^{k_{n}}
$$

- Comment 3.38. This is one advantage of the grid-based approach. We consider series only of this "multiple Laurent series" type. We do not have to contemplate series supported by abstract ordinals, something that may be considered "esoteric"-at least by beginners.
Definition 3.39. Let $\boldsymbol{\mu}=\left\{\mu_{1}, \ldots, \mu_{n}\right\} \subseteq \mathfrak{M}^{\text {small }}$ be a ratio set and $\mathbf{m} \in \mathbb{Z}^{n}$. The set of series supported by the grid $\mathfrak{J}^{\mu, \mathbf{m}}$ is

$$
\mathbb{T}^{\mu, \mathbf{m}}=\left\{T \in \mathbb{R}^{\mathfrak{M}}: \operatorname{supp} T \subseteq \mathfrak{J}^{\mu, \mathbf{m}}\right\}
$$

The set of $\boldsymbol{\mu}$-based series is

$$
\mathbb{T}^{\mu}=\bigcup_{\mathbf{m} \in \mathbb{Z}^{n}} \mathbb{T}^{\mu, \mathbf{m}}
$$

The set of grid-based series is

$$
\mathbb{R} \llbracket \mathfrak{M} \rrbracket=\bigcup_{\mu} \mathbb{T}^{\mu}
$$

In this union, all finite sets $\boldsymbol{\mu} \subset \mathfrak{M}^{\text {small }}$ are allowed, and all values of $n$ are allowed. But each individual series is supported by a grid $\mathfrak{J}^{\mu, \mathrm{m}}$ generated by one finite set $\boldsymbol{\mu}$.

If $\boldsymbol{\mu} \subseteq \widetilde{\boldsymbol{\mu}}$, then $\mathbb{T}^{\boldsymbol{\mu}} \subseteq \mathbb{T}^{\tilde{\mu}}$ in a natural way. If $\mathfrak{M}$ is a subgroup of $\widetilde{\mathfrak{M}}$ and inherits the order, then $\mathbb{R} \llbracket \mathfrak{M} \rrbracket \subseteq \mathbb{R} \llbracket \widetilde{M} \rrbracket$ in a natural way.

- Comment 3.40. The series

$$
\sum_{j=1}^{\infty} x^{1 / j}=x^{1}+x^{1 / 2}+x^{1 / 3}+x^{1 / 4}+\cdots
$$

despite having well ordered support, does not belong to $\mathbb{R} \llbracket \mathfrak{M} \rrbracket$. It is not grid-based.

- Comment 3.41. The correspondence $\mathbf{k} \mapsto \boldsymbol{\mu}^{\mathbf{k}}$ may fail to be injective. Let $\boldsymbol{\mu}=\left\{x^{-1 / 3}, x^{-1 / 2}\right\}$. Then $\mu_{1}^{3}=\mu_{2}^{2}$.
Proposition 3.42. Given $\boldsymbol{\mu}, \mathbf{m}, \mathfrak{g}$, there are only finitely many $\mathbf{k} \in \mathbf{J}_{\mathbf{m}}$ with $\boldsymbol{\mu}^{\mathbf{k}}=\mathfrak{g}$.

Proof. Suppose there are infinitely many $\mathbf{k} \in \mathbf{J}_{\mathbf{m}}$ with $\boldsymbol{\mu}^{\mathbf{k}}=\mathfrak{g}$. By Proposition 3.6, this includes $\mathbf{k}_{1}<\mathbf{k}_{2}$. But then $\boldsymbol{\mu}^{\mathbf{k}_{1}} \succ \boldsymbol{\mu}^{\mathbf{k}_{2}}$, so they are not both equal to $\mathfrak{g}$.

The map $\mathbf{k} \mapsto \boldsymbol{\mu}^{\mathbf{k}}$ might not be one-to-one, but it is finite-to-one. So: if $T_{i} \in \mathbb{T}^{\boldsymbol{\mu}, \mathbf{m}}$ for all $i \in I$, and $\mathbf{E}_{i}=\left\{\mathbf{k} \in \mathbf{J}_{\mathbf{m}}: \boldsymbol{\mu}^{\mathbf{k}} \in \operatorname{supp} T_{i}\right\}$, then $\left(\operatorname{supp} T_{i}\right)$ is point-finite if and only if $\left(\mathbf{E}_{i}\right)$ is point-finite. We may sometimes say a family $\left(T_{i}\right)$ is point-finite when the family $\left(\operatorname{supp} T_{i}\right)$ of supports is point-finite.

## Manifestly Small

Definition 3.43. If $\mathfrak{g}$ may be written in the form $\boldsymbol{\mu}^{\mathbf{k}}$ with $\mathbf{k}>\mathbf{0}$, then $\mathfrak{g}$ is $\boldsymbol{\mu}$-small, written $\mathfrak{g} \prec^{\boldsymbol{\mu}}$ 1. [For emphasis, manifestly $\boldsymbol{\mu}$-small.] If every $\mathfrak{g} \in \operatorname{supp} T$ is $\boldsymbol{\mu}$-small, then we say $T$ is $\boldsymbol{\mu}$-small, written $T \prec^{\boldsymbol{\mu}} 1$.

- Comment 3.44. Let $\boldsymbol{\mu}=\left\{x^{-1}, e^{-x}\right\}$. Then $\mathfrak{g}=x e^{-x}$ is small, but not $\boldsymbol{\mu}$-small. For $T=x^{-1}+x e^{-x}$ we have $\operatorname{mag} T \prec^{\boldsymbol{\mu}} 1$ but not $T \prec^{\boldsymbol{\mu}} 1$.


## The Asymptotic Topology

Definition 3.45. Limits of grid-based series. Let $I$ be an infinite index set (such as $\mathbb{N}$ ) and let $T_{i}, T \in \mathbb{R} \llbracket \mathfrak{M} \rrbracket$ for $i \in I$. Then: (a) $T_{i} \xrightarrow{\mu, \mathbf{m}} T$ means: $\operatorname{supp} T_{i} \subseteq \mathfrak{J}^{\boldsymbol{\mu}, \mathbf{m}}$ for all $i$, and the family $\operatorname{supp}\left(T_{i}-T\right)$ is point-finite. (b) $T_{i} \xrightarrow{\boldsymbol{\mu}}$ $T$ means: there exists $\mathbf{m}$ such that $T_{i} \xrightarrow{\boldsymbol{\mu}, \mathbf{m}} T$. (c) $T_{i} \rightarrow T$ means: there exists $\boldsymbol{\mu}$ such that $T_{i} \xrightarrow{\mu} T$. (See the "asymptotic topology" in [11, §1.2].)

- Comment 3.46. The sequence $\left(x^{j}\right)_{j \in \mathbb{N}}$ is point-finite, but it does not converge to 0 because the supports are not contained in any fixed grid $\mathfrak{J}^{\mu, \mathbf{m}}$.
- Comment 3.47. This type of convergence is not the convergence associated with the order. For example, $\left(x^{-j}\right)_{j \in \mathbb{N}} \rightarrow 0$ even though $x^{-j} \succ e^{-x}$ for all $j$. Another example: The grid-based series $\sum_{j=0}^{\infty} x^{-j}$ is $A=\left(1-x^{-1}\right)^{-1}$, even though there are many grid-based series (for example, $A-e^{-x}$ ) strictly smaller than $A$ but strictly larger than all partial sums $\sum_{j=0}^{N} x^{-j}$.

In fact, the order topology would have poor algebraic properties for sequences: For example

$$
x^{-1}>x^{-2}+e^{-x}>x^{-3}>x^{-4}+e^{-x}>\cdots
$$

(in both orderings $>$ and $\succ$ ). So in the order topology the sequences $x^{-j}$ and $x^{-j}+e^{-x}$ should have the same limit, but their difference does not converge to zero.

Proposition 3.48 (Continuity). Let $I$ be an infinite index set, and let $A_{i}, B_{i} \in$ $\mathbb{R} \llbracket \mathfrak{M} \rrbracket$ for $i \in I$. If $A_{i} \rightarrow A$ and $B_{i} \rightarrow B$, then $A_{i}+B_{i} \rightarrow A+B$ and $A_{i} B_{i} \rightarrow A B$.

Proof. We may increase $\boldsymbol{\mu}$ and decrease $\mathbf{m}$ to arrange $A_{i} \xrightarrow{\boldsymbol{\mu}, \mathbf{m}} A$ and $B_{i} \xrightarrow{\boldsymbol{\mu}, \mathbf{m}}$ $B$ for the same $\boldsymbol{\mu}, \mathbf{m}$. Then $A_{i}+B_{i} \xrightarrow{\boldsymbol{\mu}, \mathbf{m}} A+B$ and $A_{i} B_{i} \xrightarrow{\boldsymbol{\mu}, \mathbf{p}} A B$ for $\mathbf{p}=2 \mathbf{m}$. To see this: let $\mathfrak{g} \in \mathfrak{J}^{\boldsymbol{\mu}, \mathbf{p}}$. There are finitely many pairs $\mathfrak{m}, \mathfrak{n} \in \mathfrak{J}^{\boldsymbol{\mu}, \mathbf{k}}$ such that $\mathfrak{m} \mathfrak{n}=\mathfrak{g}$ (Proposition 3.27). So there is a single finite $I_{0} \subseteq I$ outside of which $A_{i}[\mathfrak{m}]=A[\mathfrak{m}]$ and $B_{i}[\mathfrak{n}]=B[\mathfrak{n}]$ for all such $\mathfrak{m}, \mathfrak{n}$. For such $i$, we also have $\left(A_{i} B_{i}\right)[\mathfrak{g}]=(A B)[\mathfrak{g}]$.

Definition 3.49. Series of grid-based series. Let $A_{i}, S \in \mathbb{R} \llbracket \mathfrak{M} \rrbracket$ for $i$ in some index set $I$. Then

$$
S=\sum_{i \in I} A_{i}
$$

means: there exist $\boldsymbol{\mu}$ and $\mathbf{m}$ such that $\operatorname{supp} A_{i} \subseteq \mathfrak{J}^{\boldsymbol{\mu}, \mathbf{m}}$ for all $i$; for all $\mathfrak{g}$, the set $I_{\mathfrak{g}}=\left\{i \in I: A_{i}[\mathfrak{g}] \neq 0\right\}$ is finite; and $S[\mathfrak{g}]=\sum_{i \in I_{\mathfrak{g}}} A_{i}[\mathfrak{g}]$.

Proposition 3.50. If $S \in \mathbb{R} \llbracket \mathfrak{M} \rrbracket$, then the "formal combination of group elements" that specifies $S$ in fact converges to $S$ in this sense as well.

Note we have the "nonarchimedean" (or "ultrametric") Cauchy criterion: In the asymptotic topology, a series $\sum A_{i}$ converges if and only if $A_{i} \rightarrow 0$.
Proposition 3.51. Let $S \in \mathbb{T}^{\mu}$ be $\boldsymbol{\mu}$-small. Then $\left(S^{j}\right)_{j \in \mathbb{N}} \xrightarrow{\mu} 0$.
Proof. Every monomial in supp $S$ can be written in the form $\boldsymbol{\mu}^{\mathbf{k}}$ with $\mathbf{k}>\mathbf{0}$. The product of two of these is again one of these. Let $\mathfrak{g}_{0} \in \mathfrak{M}$. If $\mathfrak{g}_{0}$ is not
$\boldsymbol{\mu}$-small, then $\mathfrak{g}_{0} \in \operatorname{supp}\left(S^{j}\right)$ for no $j$. So assume $\mathfrak{g}_{0}$ is $\boldsymbol{\mu}$-small. Then there are just finitely many $\mathbf{p}>\mathbf{0}$ such that $\mathfrak{g}_{0}=\boldsymbol{\mu}^{\mathbf{p}}$. Let

$$
N=\max \left\{|\mathbf{p}|: \mathbf{p}>\mathbf{0}, \boldsymbol{\mu}^{\mathbf{p}}=\mathfrak{g}_{0}\right\}
$$

Now let $j>N$. Since every $\mathfrak{g} \in \operatorname{supp} S$ is $\boldsymbol{\mu}^{\mathbf{k}}$ with $|\mathbf{k}| \geq 1$, we see that every element of $\operatorname{supp}\left(S^{j}\right)$ is $\boldsymbol{\mu}^{\mathbf{k}}$ with $|\mathbf{k}| \geq j>N$. So $\mathfrak{g}_{0} \notin \operatorname{supp}\left(S^{j}\right)$. This shows the family $\left(\operatorname{supp}\left(S^{j}\right)\right)$ is point-finite.

Proposition 3.52. Let $\boldsymbol{\mu} \subseteq \mathfrak{M}^{\text {small }}$ have $n$ elements. (a) Let $T \in \mathbb{T}^{\mu}$ be small. Then there is a (possibly larger) finite set $\widetilde{\boldsymbol{\mu}} \subseteq \mathfrak{M}^{\text {small }}$ such that $T$ is manifestly $\widetilde{\boldsymbol{\mu}}$-small. (b) Let $\mathbf{m} \in \mathbb{Z}^{n}$. There is a finite set $\widetilde{\boldsymbol{\mu}} \subseteq \mathfrak{M}^{\text {small }}$ such that $\mathfrak{J}^{\boldsymbol{\mu}, \mathbf{m}} \cap$ $\mathfrak{M}^{\text {small }} \subseteq \mathfrak{J}^{\tilde{\mu}, \mathbf{0}} \backslash\{1\}$. (See [37, Proposition 2.1], [13, 4.168], [18, Lemma 7.8].)
Proof. (a) follows from (b). Let $\mathbf{E}=\left\{\mathbf{k} \in \mathbf{J}_{\mathbf{m}}: \boldsymbol{\mu}^{\mathbf{k}} \prec 1\right\}$, so that $\mathfrak{J}^{\boldsymbol{\mu}, \mathbf{m}} \cap$ $\mathfrak{M}^{\text {small }}=\left\{\boldsymbol{\mu}^{\mathbf{k}}: \mathbf{k} \in \mathbf{E}\right\}$. By Proposition 3.7, $\operatorname{Min} \mathbf{E}$ is finite. Let $\tilde{\boldsymbol{\mu}}=\boldsymbol{\mu} \cup$ $\left\{\boldsymbol{\mu}^{\mathbf{k}}: \mathbf{k} \in \operatorname{Min} \mathbf{E}\right\}$. Note $\widetilde{\boldsymbol{\mu}} \subset \mathfrak{M}^{\text {small }}$. It is the original set $\boldsymbol{\mu}$ together with finitely many additional elements. Now for any $\mathfrak{g} \in \mathfrak{J}^{\mu, \mathbf{m}} \cap \mathfrak{M}^{\text {small }}$, there is $\mathbf{p} \in \mathbf{E}$ with $\boldsymbol{\mu}^{\mathbf{p}}=\mathfrak{g}$, and then there is $\mathbf{k} \in \operatorname{Min} \mathbf{E}$ with $\mathbf{p} \geq \mathbf{k}$, so that $\mathfrak{m}=\boldsymbol{\mu}^{\mathbf{k}} \in \widetilde{\boldsymbol{\mu}}$ and $\mathfrak{g}=\mathfrak{m} \boldsymbol{\mu}^{\mathbf{p}-\mathbf{k}}$. But $\boldsymbol{\mu}^{\mathbf{p}-\mathbf{k}}$ is $\boldsymbol{\mu}$-small and $\mathfrak{m} \in \widetilde{\boldsymbol{\mu}}$, so $\mathfrak{g}$ is manifestly $\widetilde{\boldsymbol{\mu}}$-small.

Call the set $\widetilde{\boldsymbol{\mu}} \backslash \boldsymbol{\mu}$ in (a) the smallness addendum for $T$. Sometimes I have used "smallness addendum" for $\tilde{\boldsymbol{\mu}}$ itself, but perhaps another term like "smallness extension" would be better.

- Comment 3.53. Continue Comment 3.44: If $\boldsymbol{\mu}=\left\{x^{-1}, e^{-x}\right\}$ then $x e^{-x}$ is small but not $\boldsymbol{\mu}$-small. But if we change to $\widetilde{\boldsymbol{\mu}}=\left\{x^{-1}, x e^{-x}, e^{-x}\right\}$, then $x e^{-x}$ is $\widetilde{\boldsymbol{\mu}}$-small.
- Comment 3.54. The statement like Proposition 3.52 for purely large $T$ is false. The grid-based series

$$
T=\sum_{j=0}^{\infty} x^{-j} e^{x}
$$

is purely large, but there is no finite set $\boldsymbol{\mu} \subseteq \mathfrak{M}^{\text {small }}$ and multi-index $\mathbf{m}$ such that all $x^{-j} e^{x}$ have the form $\boldsymbol{\mu}^{\mathbf{k}}$ with $\mathbf{m} \leq \mathbf{k}<\mathbf{0}$. This is because the set $\{\mathbf{k}: \mathbf{m} \leq \mathbf{k}<\mathbf{0}\}$ is finite.

Proposition 3.55. Let $S \in \mathbb{R} \llbracket \mathfrak{M} \rrbracket$ be small. Then $\left(S^{j}\right)_{j \in \mathbb{N}} \rightarrow 0$.
Proof. First, $S \in \mathbb{T}^{\mu}$ for some $\boldsymbol{\mu}$. Then $S$ is manifestly $\tilde{\boldsymbol{\mu}}$-small for some $\widetilde{\boldsymbol{\mu}} \supseteq \boldsymbol{\mu}$. Therefore $S^{j} \xrightarrow{\tilde{\mu}} 0$ by Proposition 3.51, so $S^{j} \rightarrow 0$.

Proposition 3.56. Let $\sum_{j=0}^{\infty} c_{j} z^{j}$ be a power series (not assumed to have positive radius of convergence). If $S$ is a small grid-based series, then $\sum_{j=0}^{\infty} c_{j} S^{j}$ converges in the asymptotic topology.

Proof. Use Proposition 3.52. We need to add the smallness addendum of $S$ to $\boldsymbol{\mu}$ to get a set $\widetilde{\boldsymbol{\mu}}$ such that $\sum_{j=0}^{\infty} c_{j} S^{j}$ is $\widetilde{\boldsymbol{\mu}}$-convergent.

- Comment 3.57. Continue Comment 3.53: If $\boldsymbol{\mu}=\left\{x^{-1}, e^{-x}\right\}$ then $S=x e^{-x}$ belongs to $\mathbb{T}^{\boldsymbol{\mu}}$ and is small but the series $\sum_{j=0}^{\infty} S^{j}$ is not $\boldsymbol{\mu}$-convergent. Increase to $\widetilde{\boldsymbol{\mu}}=\left\{x^{-1}, x e^{-x}, e^{-x}\right\}$ and then $\sum_{j=0}^{\infty} S^{j}$ is $\tilde{\boldsymbol{\mu}}$-convergent.
Proposition 3.58. Let $S_{1}, \ldots, S_{m}$ be $\boldsymbol{\mu}$-small grid-based series and let $p_{1}, \cdots$, $p_{m} \in \mathbb{Z}$. Then the family

$$
\left\{\operatorname{supp}\left(S_{1}^{j_{1}} S_{2}^{j_{2}} \cdots S_{m}^{j_{m}}\right): j_{1} \geq p_{1}, \ldots, j_{m} \geq p_{m}\right\}
$$

is point-finite. That is, all multiple Laurent series of the form

$$
\sum_{j_{1}=p_{1}}^{\infty} \sum_{j_{2}=p_{2}}^{\infty} \ldots \sum_{j_{m}=p_{m}}^{\infty} c_{j_{1} j_{2} \ldots j_{m}} S_{1}^{j_{1}} \cdots S_{m}^{j_{m}}
$$

are $\boldsymbol{\mu}$-convergent.
Proof. An induction on $m$ shows that we may assume $p_{1}=\cdots=p_{m}=1$, since the series with general $p_{i}$ and the series with all $p_{i}=1$, differ from each other by a finite number of series with fewer summations. So assume $p_{1}=\cdots=p_{m}=1$.

Let $\mathfrak{g}_{0} \in \mathfrak{M}$. If $\mathfrak{g}_{0}$ is not $\boldsymbol{\mu}$-small, then $\mathfrak{g}_{0} \in \operatorname{supp}\left(S_{1}^{j_{1}} \cdots S_{m}^{j_{m}}\right)$ for no $j_{1}, \ldots, j_{m}$. So assume $\mathfrak{g}_{0}$ is $\boldsymbol{\mu}$-small. There are finitely many $\mathbf{k}>\mathbf{0}$ so that $\boldsymbol{\mu}^{\mathbf{k}}=\mathfrak{g}_{0}$. Let

$$
N=\max \left\{|\mathbf{k}|: \mathbf{k}>\mathbf{0}, \boldsymbol{\mu}^{\mathbf{k}}=\mathfrak{g}_{0}\right\} .
$$

Each monomial in each $\operatorname{supp} S_{i}$ has the form $\boldsymbol{\mu}^{\mathbf{k}}$ with $|\mathbf{k}| \geq 1$. So if $j_{1}+\cdots+$ $j_{m}>N$, we have $\mathfrak{g}_{0} \notin \operatorname{supp}\left(S_{1}^{j_{1}} \cdots S_{m}^{j_{m}}\right)$.

Proposition 3.59. Let $A \in \mathbb{T}^{\boldsymbol{\mu}}$ be nonzero. Then there is a (possibly larger) finite set $\widetilde{\boldsymbol{\mu}} \subseteq \mathfrak{M}^{\text {small }}$ and $B \in \mathbb{T}^{\tilde{\mu}}$ such that $B A=1$. The set $\mathbb{R} \llbracket \mathfrak{M} \rrbracket$ of all grid-based series supported by a group $\mathfrak{M}$ is a field.

Proof. Write $A=c \boldsymbol{\mu}^{\mathbf{k}}(1+S)$, as in $3.29, \mathbf{k} \in \mathbb{Z}^{n}$. Then the inverse $B$ is:

$$
B=c^{-1} \boldsymbol{\mu}^{-\mathbf{k}} \sum_{j=0}^{\infty}(-1)^{j} S^{j}
$$

Now $c^{-1}$ is computed in $\mathbb{R}$. For the series, use Proposition 3.56. Let $\widetilde{\boldsymbol{\mu}}$ be $\boldsymbol{\mu}$ plus the smallness addendum for $S$.

We will call $\widetilde{\boldsymbol{\mu}} \backslash \boldsymbol{\mu}$ the inversion addendum for $A$.

- Comment 3.60. Continue Comment 3.57: If $\boldsymbol{\mu}=\left\{x^{-1}, e^{-x}\right\}$ and $A=$ $1+x e^{-x}$, then $A \in \mathbb{T}^{\mu}$. But $A$ has no inverse in $\mathbb{T}^{\mu}$. Increase to $\widetilde{\boldsymbol{\mu}}=$ $\left\{x^{-1}, x e^{-x}, e^{-x}\right\}$ and then $A^{-1} \in \mathbb{T}^{\tilde{\mu}}$.

The algebra $\mathbb{R} \llbracket \mathfrak{M} \rrbracket$ is an ordered field: If $S, T>0$, then $S+T>0$ and $S T>0$. Also: if $T_{i}>0$ and $\sum T_{i}$ exists, then $\sum T_{i}>0$.

- Comment 3.61. But: if $T_{i} \geq 0, T_{i} \rightarrow T$, then $T \geq 0$ need not follow. Take $T_{i}=x^{-i} e^{x}-x$ and $T=-x$. Also: $S, T \succ 1$ need not imply $S+T \succ 1$. For example, $S=x, T=-x+e^{-x}$.


## 3D Transseries for $x \rightarrow+\infty$.

For (real, grid-based) transseries, we define a specific ordered group $\mathfrak{G}$ of transmonomials to use for $\mathfrak{M}$. This is done in stages.

- Comment 3.62. A symbol " $x$ " appears in the notation. When we think of a transseries as describing behavior as $x \rightarrow+\infty$, then $x$ is supposed to be a large parameter. When we write "compositions" involving transseries, $x$ represents the identity function. But usually it is just a convenient symbol.
Notation 3.63. Group $\mathfrak{G}_{0}$ is isomorphic to $\mathbb{R}$ with addition and the usual ordering. To fit our applications, we write $x^{b}$ for the group element corresponding to $b \in \mathbb{R}$. Then $x^{a} x^{b}=x^{a+b} ; x^{0}=1 ; x^{-b}$ is the inverse of $x^{b} ; x^{a} \prec x^{b}$ iff $a<b$.

Log-free transseries of height zero are those obtained from this group as in Definition 3.39. Write $\mathbb{T}_{0}=\mathbb{R} \llbracket \mathfrak{G}_{0} \rrbracket$. Then the set of purely large transseries in $\mathbb{T}_{0}$ (including 0 ) is a group under addition.

- Comment 3.64. Transseries of height zero:

$$
-x^{3}+2 x^{2}-x, \quad \sum_{j=1}^{\infty} \sum_{k=1}^{\infty} x^{-j-k \sqrt{2}}
$$

The first is purely large, the second is small.
$\circ$
Notation 3.65. Group $\mathfrak{G}_{1}$ consists of ordered pairs $(b, L)$ —but written $x^{b} e^{L}$ where $b \in \mathbb{R}$ and $L \in \mathbb{T}_{0}$ is purely large. Define the group operation: $\left(x^{b_{1}} e^{L_{1}}\right)\left(x^{b_{2}} e^{L_{2}}\right)=x^{b_{1}+b_{2}} e^{L_{1}+L_{2}}$. Define order lexicographically: $\left(x^{b_{1}} e^{L_{1}}\right) \succ$ $\left(x^{b_{2}} e^{L_{2}}\right)$ iff either $L_{1}>L_{2}$ or $\left\{L_{1}=L_{2}\right.$ and $\left.b_{1}>b_{2}\right\}$. Identify $\mathfrak{G}_{0}$ as a subgroup of $\mathfrak{G}_{1}$, where $x^{b}$ is identified with $x^{b} e^{0}$.

Log-free transseries of height 1 are those obtained from this group as in Definition 3.39. Write $\mathbb{T}_{1}=\mathbb{R} \llbracket \mathfrak{G}_{1} \rrbracket$. We may identify $\mathbb{T}_{0}$ as a subset of $\mathbb{T}_{1}$. Then the set of purely large transseries in $\mathbb{T}_{1}$ (including 0 ) is a group under addition.

- Comment 3.66. Transseries of height 1:

$$
e^{-x^{3}+2 x^{2}-x}, \quad \sum_{j=1}^{\infty} x^{-j} e^{x}, \quad x^{3}+e^{-x^{3 / 4}}
$$

The first is small, the second is purely large, the last is large but not purely large.
Notation 3.67. Suppose log-free transmonomials $\mathfrak{G}_{N}$ and log-free transseries $\mathbb{T}_{N}$ of height $N$ have been defined. Group $\mathfrak{G}_{N+1}$ consists of ordered pairs ( $b, L$ ) but written $x^{b} e^{L}$, where $b \in \mathbb{R}$ and $L \in \mathbb{T}_{N}$ is purely large. Define the group operation: $\left(x^{b_{1}} e^{L_{1}}\right)\left(x^{b_{2}} e^{L_{2}}\right)=x^{b_{1}+b_{2}} e^{L_{1}+L_{2}}$. Define order: $\left(x^{b_{1}} e^{L_{1}}\right) \succ$ $\left(x^{b_{2}} e^{L_{2}}\right)$ iff either $L_{1}>L_{2}$ or $\left\{L_{1}=L_{2}\right.$ and $\left.b_{1}>b_{2}\right\}$.

Identify $\mathfrak{G}_{N}$ as a subgroup of $\mathfrak{G}_{N+1}$ recursively.
Log-free transseries of height $N+1$ are those obtained from this group as in Definition 3.39. Write $\mathbb{T}_{N+1}=\mathbb{R} \llbracket \mathfrak{G}_{N+1} \rrbracket$. We may identify $\mathbb{T}_{N}$ as a subset of $\mathbb{T}_{N+1}$.

- Comment 3.68. Height 2: $e^{-e^{x}}, e^{\sum_{j=1}^{\infty} x^{-j} e^{x}}$.

Notation 3.69. The group of log-free transmonomials is

$$
\mathfrak{G}_{\bullet}=\bigcup_{N \in \mathbb{N}} \mathfrak{G}_{N}
$$

The field of log-free transseries is

$$
\mathbb{T}_{\bullet}=\bigcup_{N \in \mathbb{N}} \mathbb{T}_{N}
$$

In fact, $\left.\mathbb{T}_{\bullet}=\mathbb{R} \llbracket \mathfrak{G}_{\bullet}\right]$ because each individual transseries is grid-based. Any grid in $\mathfrak{G}_{\bullet}$ is contained in $\mathfrak{G}_{N}$ for some $N$.

A ratio set $\boldsymbol{\mu}$ is hereditary if for every transmonomial $x^{b} e^{L}$ in $\boldsymbol{\mu}$, we also have $L \in \mathbb{T}^{\mu}$. Of course, given any ratio set $\boldsymbol{\mu} \subseteq \mathfrak{G}^{\text {small }}$, there is a hereditary ratio set $\widetilde{\boldsymbol{\mu}} \supseteq \boldsymbol{\mu}$. (When we add a set of generating ratios for the exponents $L$, then sets of generating ratios for their exponents, and so on, the process ends in finitely many steps, by induction on the heights. We do need to know: If $T \in \mathbb{T}_{\bullet}$, then the union

$$
\bigcup_{x^{b} e_{L} L \text { supp } T} \operatorname{supp} L
$$

is a subgrid. See [22, Prop 2.21].) Call $\widetilde{\boldsymbol{\mu}} \backslash \boldsymbol{\mu}$ the heredity addendum of $\boldsymbol{\mu}$.

Remark 3.70. If $\mathfrak{M}$ is a group, then $\mathbb{R} \llbracket \mathfrak{M}]$ is a field. In particular, $\mathbb{T}_{N}=$ $\mathbb{R} \llbracket \mathfrak{G}_{N} \rrbracket$ is a field $(N=0,1,2, \cdots)$.

Proposition 3.71. Let $A$ be a nonzero log-free transseries. If $A \succ 1$, then there exists a real number $c>0$ such that $A \succ x^{c}$. If $A \prec 1$, then there exists a real number $c<0$ such that $A \prec x^{c}$.

Proof. Let $\operatorname{mag} A=x^{b} e^{L} \succ 1$. If $L=0$, then $b>0$, so take $c=b / 2$. If $L>0, A \succ x^{1}$, since $\succ$ is defined lexicographically. The other case is similar.

Proposition 3.72 (Height Wins). Let $L>0$ be purely large of height $N$ and not $N-1$, let $b \in \mathbb{R}$, and let $T \neq 0$ be of height $N$. Then $x^{b} e^{L} \succ T$ and $x^{b} e^{-L} \prec T$.

Proof. By induction on the height. Let $\operatorname{mag} T=x^{b_{1}} e^{L_{1}}$. So $L_{1} \in \mathbb{T}_{N-1}$, and therefore by the induction hypothesis $\operatorname{dom}\left(L-L_{1}\right)=\operatorname{dom}(L)>0$. So $L>L_{1}$ and $x^{b} e^{L} \succ x^{b_{1}} e^{L_{1}}$.

If $\mathfrak{n} \in \mathfrak{G}_{N} \backslash \mathfrak{G}_{N-1}$ (we say $\mathfrak{n}$ has $\boldsymbol{e x a c t}$ height $N$ ), then either (i) $\mathfrak{n} \succ 1$ and $\mathfrak{n} \succ \mathfrak{m}$ for all $\mathfrak{m} \in \mathfrak{G}_{N-1}$, or (ii) $\mathfrak{n} \prec 1$ and $\mathfrak{n} \prec \mathfrak{m}$ for all $\mathfrak{m} \in \mathfrak{G}_{N-1}$. [We say $\mathfrak{G}_{N-1}$ is convex in $\left.\mathfrak{G}_{N}.\right]$

- Comment 3.73. $\mathfrak{n}=e^{-e^{x}}$ has exact height 2, and $T=\sum_{k=1}^{\infty} \sum_{j=1}^{\infty} x^{-j} e^{-k x}$ has height 1 , so of course $\mathfrak{n} \prec T$. Even more: $T / \mathfrak{n}$ is purely large.


## Derivative

Definition 3.74. Derivative (notations ', $\partial$ ) is defined recursively. First, $\left(x^{a}\right)^{\prime}=a x^{a-1}$. (If we are keeping track of generating ratios, we may need the addendum of ratio $x^{-1}$.) If $\partial$ has been defined for $\mathfrak{G}_{N}$, then define it termwise for $\mathbb{T}_{N}$ :

$$
\begin{equation*}
\left(\sum a_{\mathfrak{g}} \mathfrak{g}\right)^{\prime}=\sum a_{\mathfrak{g}} \mathfrak{g}^{\prime} \tag{*}
\end{equation*}
$$

(See the next proposition for the proof that this makes sense.) Then, if $\partial$ has been defined for $\mathbb{T}_{N}$, define it on $\mathfrak{G}_{N+1}$ by

$$
\left(x^{b} e^{L}\right)^{\prime}=b x^{b-1} e^{L}+x^{b} L^{\prime} e^{L}=\left(b x^{-1}+L^{\prime}\right) x^{b} e^{L}
$$

For the derivative addendum $\widetilde{\boldsymbol{\mu}} \backslash \boldsymbol{\mu}$ : begin with $\boldsymbol{\mu}$, add the heredity addendum of $\boldsymbol{\mu}$, and add $x^{-1}$. So (by induction) if $T \in \mathbb{T}^{\mu}$, then $T^{\prime} \in \mathbb{T}^{\tilde{\mu}}$. Repeating the derivative addendum adds nothing new, so in fact all derivatives $T^{(j)}$ belong to $\mathbb{T}^{\tilde{\mu}}$.

Remark 3.75. This derivative satisfies all the usual algebraic properties of the derivative. There are just lots of tedious things to check. $(A B)^{\prime}=A^{\prime} B+A B^{\prime}$, $\left(S^{k}\right)^{\prime}=k S^{\prime} S^{k-1}$, etc.

Here is the computation needed for the induction step when we recursively define the derivative using $(*)$.

Proposition 3.76. Let $\boldsymbol{\mu}$ be given. Let $\widetilde{\boldsymbol{\mu}}$ be as described. (i) If $T_{i} \xrightarrow{\boldsymbol{\mu}} T$ then $T_{i}^{\prime} \xrightarrow{\tilde{\boldsymbol{\mu}}} T^{\prime}$. (ii) If $\sum T_{i}$ is $\boldsymbol{\mu}$-convergent, then $\sum T_{i}^{\prime}$ is $\tilde{\boldsymbol{\mu}}$-convergent and $\left(\sum T_{i}\right)^{\prime}=\sum T_{i}^{\prime}$. (iii) If $\mathfrak{A} \subseteq \mathfrak{J}^{\boldsymbol{\mu}, \mathbf{m}}$, then $\sum_{\mathfrak{g} \in \mathfrak{A}} a_{\mathfrak{g}} \mathfrak{g}^{\prime}$ is $\widetilde{\boldsymbol{\mu}}$-convergent.

Proof. (iii) is stated equivalently: the family (supp $\mathfrak{g}^{\prime}$ ) is point-finite. Or: as $\mathfrak{g}$ ranges over $\mathfrak{f}^{\boldsymbol{\mu}, \mathbf{m}}$, we have $\mathfrak{g}^{\prime} \xrightarrow{\tilde{\mu}} 0$.

Proof by induction on the height.
Say $\mu_{1}=x^{-b_{1}} e^{-L_{1}}, \ldots, \mu_{n}=x^{-b_{n}} e^{-L_{n}}$, and $\mathbf{k}=\left(k_{1}, \ldots, k_{n}\right)$. Then

$$
\begin{aligned}
\left(\boldsymbol{\mu}^{\mathbf{k}}\right)^{\prime} & =\left(x^{-k_{1} b_{1}-\cdots-k_{n} b_{n}} e^{-k_{1} L_{1}-\cdots-k_{n} L_{n}}\right)^{\prime} \\
& =\left(-k_{1} b_{1}-\cdots-k_{n} b_{n}\right) x^{-1} \boldsymbol{\mu}^{\mathbf{k}}+\left(-k_{1} L_{1}^{\prime}-\cdots-k_{n} L_{n}^{\prime}\right) \boldsymbol{\mu}^{\mathbf{k}}
\end{aligned}
$$

So if $T=\sum_{\mathbf{k} \geq \mathbf{m}} a_{\mathbf{k}} \boldsymbol{\mu}^{\mathbf{k}}$, then summing the above transmonomial result, we get

$$
T^{\prime}=x^{-1} T_{0}+L_{1}^{\prime} T_{1}+\cdots+L_{n}^{\prime} T_{n}
$$

where $T_{0}, \ldots, T_{n}$ are transseries with the same support as $T$, and therefore they exist in $\mathbb{T}^{\mu, \mathbf{m}}$. Derivatives $L_{1}^{\prime}, \ldots, L_{n}^{\prime}$ exist by induction hypothesis. So $T^{\prime}$ exists.

The preceding proof suggest the following. We introduce lsupp, thought of as "the support of the logarithmic derivative" for monomials.
Definition 3.77. For (log-free) monomials, define

$$
\operatorname{lsupp}\left(x^{b} e^{L}\right)=\left\{x^{-1}\right\} \cup \operatorname{supp} L^{\prime}
$$

For a set $\mathfrak{A} \subseteq \mathfrak{G}_{\bullet}$, define $\operatorname{lsupp} \mathfrak{A}=\bigcup_{\mathfrak{g} \in \mathfrak{A}} \operatorname{lsupp} \mathfrak{g}$. For $T \in \mathbb{T}_{\bullet}$, define lsupp $T=$ $\operatorname{lsupp} \operatorname{supp} T$. For a set $\mathbf{A} \subseteq \mathbb{T}_{\bullet}$, define $\operatorname{lsupp} \mathbf{A}=\bigcup_{T \in \mathbf{A}} \operatorname{lsupp} T$.

Proposition 3.78. Properties of lsupp.
(a) If $\boldsymbol{\mu}=\left\{\mu_{1}, \ldots, \mu_{n}\right\} \subseteq \mathfrak{G}_{\bullet}^{\text {small }}$ and $\mathbf{k} \in \mathbb{Z}^{n}$, then $\operatorname{lsupp} \boldsymbol{\mu}^{\mathbf{k}} \subseteq \operatorname{lsupp} \boldsymbol{\mu}$. So $\operatorname{lsupp} \mathbb{T}^{\mu} \subseteq \operatorname{lsupp} \boldsymbol{\mu}$.
(b) For any finite $\boldsymbol{\mu} \subseteq \mathfrak{G}_{\bullet}^{\text {small }}, \operatorname{lsupp}(\boldsymbol{\mu})$ is a subgrid.
(c) If $T \in \mathbb{T}^{\boldsymbol{\mu}}$, then $\operatorname{supp}\left(T^{\prime}\right) \subseteq \operatorname{lsupp}(\boldsymbol{\mu}) \cdot \operatorname{supp}(T)$.
(d) If $\widetilde{\boldsymbol{\mu}}$ is the smallness addendum for some $S \in \mathbb{T}^{\boldsymbol{\mu}}$, then $\operatorname{lsupp} \widetilde{\boldsymbol{\mu}} \subseteq$ $\operatorname{lsupp} \boldsymbol{\mu}$.
(e) If $\mathfrak{g} \in \mathfrak{G}_{N}, N \geq 1$, then $\operatorname{lsupp} \mathfrak{g} \subseteq \mathfrak{G}_{N-1}$.
(f) If $\mathfrak{g} \in \mathfrak{G}_{0}, \mathfrak{g} \neq 1$, then $\operatorname{lsupp} \mathfrak{g}=\left\{x^{-1}\right\}$.
(g) Each $\mathbb{T}_{N}$ is closed under $\partial$.

Proof. (a) Say $\mu_{1}=x^{-b_{1}} e^{-L_{1}}, \ldots, \mu_{n}=x^{-b_{n}} e^{-L_{n}}$. Then

$$
\begin{aligned}
\boldsymbol{\mu}^{\mathbf{k}} & =x^{-k_{1} b_{1}-\cdots-k_{n} b_{n}} e^{-k_{1} L_{1}-\cdots-k_{n} L_{n}}, \text { so } \\
\operatorname{lsupp} \boldsymbol{\mu}^{\mathbf{k}} & =\left\{x^{-1}\right\} \cup \operatorname{supp}\left(-k_{1} L_{1}^{\prime}-\cdots-k_{n} L_{n}^{\prime}\right) \\
& \subseteq\left\{x^{-1}\right\} \cup \operatorname{supp}\left(L_{1}^{\prime}\right) \cup \cdots \cup \operatorname{supp}\left(L_{n}^{\prime}\right)=\operatorname{lsupp} \boldsymbol{\mu} .
\end{aligned}
$$

(b) Each $\operatorname{supp} L_{i}^{\prime} \in \mathbf{W}$ and $\left\{x^{-1}\right\} \in \mathbf{W}$, so their (finite) union also belongs to $\mathbf{W}$.
(c) Proposition 3.76.
(d) Use the proof of Proposition 3.52 together with (a).
(e) and (f) are clear.
(g) Use (c) and (e).

Remark 3.79. Similar to (d): If $\widetilde{\boldsymbol{\mu}}$ is the inversion addendum for some $A \in \mathbb{T}^{\mu}$, then $\operatorname{lsupp} \widetilde{\boldsymbol{\mu}} \subseteq \operatorname{lsupp} \boldsymbol{\mu}$.

Here are a few technical results on derivatives. They lead to Proposition 3.85 , where we will prove that $T^{\prime}=0$ only if $T$ is "constant" in the sense used here, that $T \in \mathbb{R}$.

Proposition 3.80. There is no $T \in \mathbb{T}_{\bullet}$ with $T^{\prime}=x^{-1}$.
Proof. In fact, we show: If $\mathfrak{g} \in \mathfrak{G}_{\bullet}$, then $x^{-1} \notin \operatorname{supp} \mathfrak{g}^{\prime}$. This suffices since

$$
\operatorname{supp} T^{\prime} \subseteq \bigcup_{\mathfrak{g} \in \operatorname{supp} T} \operatorname{supp} \mathfrak{g}^{\prime}
$$

Proof by induction on the height. If $\mathfrak{g}=x^{b}$ has height 0 , then $\mathfrak{g}^{\prime}=b x^{b-1}$ and $x^{-1} \notin \operatorname{supp} \mathfrak{g}^{\prime}$. If $\mathfrak{g}=x^{b} e^{L}$ has exact height $N$, so $L$ is purely large of exact height $N-1$, then $\mathfrak{g}^{\prime}=\left(b x^{-1}+L^{\prime}\right) x^{b} e^{L}$. Now by the induction hypothesis, $b x^{-1}+L^{\prime} \neq 0$, so (by Proposition 3.72) every term of $\mathfrak{g}^{\prime}$ is far larger than $x^{-1}$ if $L>0$ and far smaller than $x^{-1}$ if $L<0$. So $x^{-1} \notin \operatorname{supp} \mathfrak{g}^{\prime}$.

Proposition 3.81. (a) Let $\mathfrak{m} \neq 1$ be a log-free monomial with exact height n. Then $\mathfrak{m}^{\prime}$ also has exact height $n$. (b) If $\mathfrak{m} \succ \mathfrak{n}$, and $\mathfrak{m} \neq 1$, then $\mathfrak{m}^{\prime} \succ \mathfrak{n}^{\prime}$. (c) If $\operatorname{mag} T \neq 1$, then $T^{\prime} \asymp(\operatorname{mag} T)^{\prime}$ and $T^{\prime} \sim(\operatorname{dom} T)^{\prime}$. (d) If $\operatorname{mag} T \neq 1$ and $T \succ S$, then $T^{\prime} \succ S^{\prime}$.

Proof. (a) For height $0, \mathfrak{m}=x^{b}, b \neq 0$ so $\mathfrak{m}^{\prime}=b x^{b-1} \neq 0$ also has height 0 . Let $\mathfrak{m}=x^{b} e^{L} \neq 1$ have exact height $n$, so that $L \neq 0$ has exact height $n-1$. Of course $L^{\prime}$ has height at most $n-1$, and $\left(b x^{-1}+L^{\prime}\right)$ is not zero by Proposition 3.80, so $\mathfrak{m}^{\prime}=\left(b x^{-1}+L^{\prime}\right) x^{b} e^{L}$ again has exact height $n$.

If (b) holds for all $\mathfrak{m}, \mathfrak{n}$ of a given height $n$, then (c) and (d) follow for $S, T$ of height $n$. So it remains to prove (b).

First suppose $\mathfrak{m}, \mathfrak{n}$ have different heights. Say $\mathfrak{m} \in \mathfrak{G}_{m}, \mathfrak{n} \in \mathfrak{G}_{n} \backslash \mathfrak{G}_{m}, n>m$. If $\mathfrak{n} \succ 1$, then $\mathfrak{n}^{\prime} \succ 1$ (because its height $n>0$ ), $\mathfrak{n}^{\prime} \in \mathfrak{G}_{n} \backslash \mathfrak{G}_{m}, \mathfrak{m} \prec \mathfrak{n}$ by Proposition 3.72, $\mathfrak{m}^{\prime} \prec \mathfrak{n}^{\prime}$. If $\mathfrak{n} \prec 1$, then $\mathfrak{n}^{\prime} \prec 1, \mathfrak{n}^{\prime} \in \mathfrak{G}_{n} \backslash \mathfrak{G}_{m}, \mathfrak{m} \succ \mathfrak{n}, \mathfrak{m}^{\prime} \succ \mathfrak{n}^{\prime}$. So (b) holds in both cases. That is: either $\left\{\mathfrak{m} \succ \mathfrak{n}\right.$ and $\left.\mathfrak{m}^{\prime} \succ \mathfrak{n}^{\prime}\right\}$ or $\{\mathfrak{m} \prec \mathfrak{n}$ and $\left.\mathfrak{m}^{\prime} \prec \mathfrak{n}^{\prime}\right\}$.

So suppose $\mathfrak{m}, \mathfrak{n}$ have the same height $n$. Write $\mathfrak{m}=x^{a} e^{A}, \mathfrak{n}=x^{b} e^{B}$, where $a, b$ real and $A, B$ purely large. Assume $\mathfrak{m} \succ \mathfrak{n}$, so either $A>B$ or $A=B, a>b$. We take the case $A>B$ (the other one is similar to Case 2, below). Say $A-B$ has exact height $k$. There will be two cases: $k=n-1$ and $k<n-1$.

Case 1. $k=n-1$. Then $A-B$ has exact height $n-1$ and

$$
\frac{b x^{-1}+B^{\prime}}{a x^{-1}+A^{\prime}}
$$

has height $n-1$ (and its denominator is not zero by Proposition 3.80). Therefore by Proposition 3.72,

$$
x^{a-b} e^{A-B} \succ \frac{b x^{-1}+B^{\prime}}{a x^{-1}+A^{\prime}},
$$

and thus $\left(a x^{-1}+A^{\prime}\right) x^{a} e^{A} \succ\left(b x^{-1}+B^{\prime}\right) x^{b} e^{B}$. That is, $\mathfrak{m}^{\prime} \succ \mathfrak{n}^{\prime}$.
Case 2. $k<n-1$. Write $A=A_{0}+A_{1}, B=B_{0}+A_{1}$ where purely large $A_{0}, B_{0}$ have height $k$ (and purely large $A_{1} \neq 0$ has exact height $n-1>k$ ). Now $A_{1}^{\prime}$ has height $n-1>k$ and is large, so $A_{1}^{\prime} \succ a x^{-1}+A_{0}^{\prime}$ and $A_{1}^{\prime} \succ b x^{-1}+B_{0}^{\prime}$. Since $x^{a-b} e^{A_{0}-B_{0}} \succ 1$, we have $x^{a} e^{A_{0}} \succ x^{b} e^{B_{0}}$ and therefore

$$
\begin{aligned}
& \mathfrak{m}^{\prime}=\left(a x^{-1}+A_{0}^{\prime}+A_{1}^{\prime}\right) x^{a} e^{A_{0}+A_{1}} \sim A_{1}^{\prime} x^{a} e^{A_{0}+A_{1}} \\
& \succ A_{1}^{\prime} x^{b} e^{B_{0}+A_{1}} \sim\left(b x^{-1}+B_{0}^{\prime}+A_{1}^{\prime}\right) x^{b} e^{B_{0}+A_{1}}=\mathfrak{n}^{\prime}
\end{aligned}
$$

(See [18, Prop. 4.1] for this proof.)

Proposition 3.82. Let $T \in \mathbb{T}_{\bullet}$. (i) If $T \prec 1$, then $T^{\prime} \prec 1$. (ii) If $T \succ 1$ and $T>0$, then $T^{\prime}>0$. (iii) If $T \succ 1$ and $T<0$, then $T^{\prime}<0$. (iv) If $T \succ 1$, then $T^{2} \succ T^{\prime}$. (v) If $T \succ 1$, then $\left(T^{\prime}\right)^{2} \succ T^{\prime \prime}$. (vi) If $T \succ 1$ then $x T^{\prime} \succ 1$.

Proof. (i) $T \prec 1 \quad \Longrightarrow \quad T \prec x \quad \Longrightarrow \quad T^{\prime} \prec 1$.
(ii)(iii) Assume $T \succ 1$. Let $\operatorname{dom} T=a x^{b} e^{L}$, so $T$ has the same sign as $a$. Then $T^{\prime} \sim a\left(b x^{-1}+L^{\prime}\right) x^{b} e^{L}$. The proof is by induction on the height of $T$. If $T$ has height 0 , so that $L=0$ and $b>0$, then $T^{\prime} \sim a b x^{b-1} e^{L}$ has the same sign as $a$. Assume $T$ has height $N>0$, so $L>0$ and $L$ has height $N-1$, so (since $L$ is large) the induction hypothesis tells us that $L^{\prime}>0$. Also, $L \succ x^{c}$ for some $c>0$ so $L^{\prime} \succ x^{c-1} \succ x^{-1}$, so $T^{\prime} \sim a b L^{\prime} x^{b} e^{L}$ has the same sign as $a$.
(iv) $T \succ 1 \Longrightarrow T \succ \frac{1}{x} \Longrightarrow \frac{1}{T} \prec x \Longrightarrow\left(\frac{1}{T}\right)^{\prime} \prec 1 \Longrightarrow \frac{T^{\prime}}{T^{2}} \prec 1$.
(v) $T \succ x^{c}$ for some $c>0$, so $T^{\prime} \succ x^{c-1} \succ x^{-1}$, then proceed as in (iv).
(vi) $T^{\prime} \succ 1 / x$ as in the proof of (v).

- Comment 3.83. After we define real powers (Definition 3.87), we will be able to formulate variants of (iv): If $T \succcurlyeq 1$, then $|T|^{1+\varepsilon} \succ T^{\prime}$ for all real $\varepsilon>0$; if $T \prec 1$, then $|T|^{1-\varepsilon} \succ T^{\prime}$ for all real $\varepsilon>0$ [18, Prop. 4.1]. And if $T>x^{a}$ for all real $a$, then $T^{1-\varepsilon} \prec T^{\prime} \prec T^{1+\varepsilon}$ for all real $\varepsilon>0$ [18, Cor. 4.4]. The proofs are essentially as given above for (iv).

Proposition 3.84. (a) If $L \neq 0$ is large and $b \in \mathbb{R}$, then $\operatorname{dom}\left(\left(x^{b} e^{L}\right)^{\prime}\right)=$ $x^{b} e^{L} \operatorname{dom}\left(L^{\prime}\right)$. (b) If $\mathfrak{g} \in \mathfrak{G}_{\bullet}, \mathfrak{g} \neq 1$, then $\mathfrak{g}^{\prime} \neq 0$.

Proof. (a) Since $L \succ 1$, there is $c>0$ with $L \succ x^{c}$, so $L^{\prime} \succ x^{c-1} \succ x^{-1}$. So $\left(x^{b} e^{L}\right)^{\prime}=x^{b} e^{L}\left(b x^{-1}+L^{\prime}\right) \asymp x^{b} e^{L} L^{\prime}$. For (b), use induction on the height and (a).

Proposition 3.85. Let $T \in \mathbb{T}_{\bullet}$. If $T^{\prime}=0$, then $T$ is a constant.
Proof. Assume $T^{\prime}=0$. Write $T=L+c+S$ as in 3.24. If $L \neq 0$ then $T^{\prime} \asymp(\operatorname{mag} L)^{\prime} \neq 0$ so $T^{\prime} \neq 0$. If $L=0$ and $S \neq 0$, then $T^{\prime} \asymp(\operatorname{mag} S)^{\prime} \neq 0$ so $T^{\prime} \neq 0$. Therefore $T=c$.

The set $\mathbb{T}_{N}$ is a differential field with constants $\mathbb{R}$. This means it follows the rules you already know for computations involving derivatives.

Proposition 3.86 (Addendum Height). Let $\boldsymbol{\mu} \subseteq \mathfrak{G}_{N}^{\text {small }}$, and let $T \in \mathbb{T}^{\boldsymbol{\mu}}$. (i) If $\widetilde{\boldsymbol{\mu}}$ is the smallness addendum for $T$, then $\widetilde{\boldsymbol{\mu}} \subseteq \mathfrak{G}_{N}^{\text {small }}$. (ii) If $\widetilde{\boldsymbol{\mu}}$ is the inversion addendum for $T$, then $\widetilde{\boldsymbol{\mu}} \subseteq \mathfrak{G}_{N}^{\text {small }}$. (iii) If $\widetilde{\boldsymbol{\mu}}$ is the heredity addendum for $\boldsymbol{\mu}$, then $\widetilde{\boldsymbol{\mu}} \subseteq \mathfrak{G}_{N}^{\text {small }}$. (iv) If $\widetilde{\boldsymbol{\mu}}$ is the derivative addendum for $T$, then $\widetilde{\boldsymbol{\mu}} \subseteq \mathfrak{G}_{N}^{\text {small }}$.

## Compositions

The field of transseries has an operation of "composition." The result $T \circ S$ is, however, defined in general only for some $S$. We will start with the easy cases.
Definition 3.87. We define $S^{b}$, where $S \in \mathbb{T}_{\bullet}$ is positive, and $b \in \mathbb{R}$. First, write $S=c x^{a} e^{L}(1+U)$ as in 3.29, with $c>0$. Then define $S^{b}=c^{b} x^{a b} e^{b L}(1+U)^{b}$. Constant $c^{b}$, with $c>0$, is computed in $\mathbb{R}$. Next, $x^{a b}$ is a transmonomial, but (if we are keeping track of generating ratios) may require addendum of a ratio. Also, $(1+U)^{b}$ is a convergent binomial series, again we may require the smallness addendum for $U$. Finally, since $L$ is purely large, so is $b L$, and thus $e^{b L}$ is a transmonomial, but may require addendum of a ratio. Of course, when $b$ is a positive integer, then the new meaning of $S^{b}$ agrees with its former meaning in terms of repeated multiplication.
Remark 3.88. Note $S^{b}$ is not of greater height than $S$ : If $S \in \mathbb{T}_{N}$, then $S^{b} \in \mathbb{T}_{N}$. If $b \neq 0$, then because $\left(S^{b}\right)^{1 / b}=S$, in fact the exact height of $S$ is the same as $S^{b}$.

- Comment 3.89. Monotonicity: If $b>0$ and $S_{1}<S_{2}$, then $S_{1}^{b}<S_{2}^{b}$. If $b<0$ and $S_{1}<S_{2}$, then $S_{1}^{b}>S_{2}^{b}$.
Definition 3.90. We define $e^{S}$, where $S \in \mathbb{T}_{\bullet}$. Write $S=L+c+U$ as in 3.24, with $L$ purely large, $c$ a constant, and $U$ small. Then $e^{S}=e^{L} e^{c} e^{U}$. Constant $e^{c}$ is computed in $\mathbb{R}$. [Note that $e^{S}>0$ since the leading coefficient is $e^{c}$.] Next, $e^{U}$ is a power series (with point-finite convergence); we may need the smallness addendum for $U$. And of course $e^{L}$ is a transmonomial, but might not already be a generating ratio, so perhaps $e^{L}$ or $e^{-L}$ is required as addendum.
Remark 3.91. Of course, if $S=L$ is purely large, then this definition of $e^{S}$ agrees with the formal notation $e^{L}$ used before. Height increases by at most one: If $S \in \mathbb{T}_{N}$, then $e^{S} \in \mathbb{T}_{N+1}$.
- Comment 3.92. Monotonicity: If $S_{1}<S_{2}$ then $e^{S_{1}}<e^{S_{2}}$.

Definition 3.93 . Let $S, T \in \mathbb{T}_{\bullet}$ with $S$ positive and large (but not necessarily purely large). We want to define the composition $T \circ S$. This is done by induction on the height of $T$. When $T=x^{b} e^{L}$ is a transmonomial, define $T \circ S=S^{b} e^{L \circ S}$. Both $S^{b}$ and $e^{L \circ S}$ may require addenda. And $L \circ S$ exists by the induction hypothesis. In general, when $T=\sum c_{\mathfrak{g}} \mathfrak{g}$, define

$$
\begin{equation*}
T \circ S=\sum c_{\mathfrak{g}}(\mathfrak{g} \circ S) \tag{*}
\end{equation*}
$$

The next proposition is required.

Remark 3.94. If $T \succ 1$, then $T \circ S \succ 1$. If $T \prec 1$, then $T \circ S \prec 1$. Because of our use of the symbol $x$, it will not be unexpected if we sometimes write $T(S)$ for $T \circ S$. Alternate term: "large and positive" = "infinitely increasing".

Here is the computation needed for the induction step when we recursively define composition using (*).

Proposition 3.95. Let $\mathfrak{J}^{\mu, \mathbf{m}}$ be a log-free grid and let $S \in \mathbb{T} \bullet$ be a large, positive, log-free transseries. Then there exist $\tilde{\boldsymbol{\mu}}$ and $\widetilde{\mathbf{m}}$ so that $\mathfrak{g} \circ S \in \mathbb{T}^{\tilde{\mu}, \tilde{\mathbf{m}}}$ for all $\mathfrak{g} \in \mathfrak{J}^{\mu, \mathbf{m}}$, and the family $(\operatorname{supp}(\mathfrak{g} \circ S))$ is point-finite.
Proof. First, add the heredity addendum of $\boldsymbol{\mu}$. Now for the ratios $\mu_{1}, \ldots, \mu_{m}$, write $\mu_{i}=x^{-b_{i}} e^{-L_{i}}, 1 \leq i \leq m$. Arrange the list so that for all $i, L_{i} \in$ $\mathbb{T}^{\left\{\mu_{1}, \ldots, \mu_{i-1}\right\}}$. Then take the $\mu_{i}$ in order. Each $S^{-b_{i}}$ may require an addendum. Each $L_{i} \circ S$ may require an addendum. So all $\mu_{i} \circ S$ exist. They are small. Add smallness addenda for these. So finally we get $\widetilde{\boldsymbol{\mu}}$.

Now for each $\mu_{i} \in \boldsymbol{\mu}$, we have $\mu_{i} \circ S$ is $\widetilde{\boldsymbol{\mu}}$-small. So by Proposition 3.58 we have $(\mathfrak{g} \circ S)_{\mathfrak{g} \in \mathfrak{J}^{\mu, \mathrm{m}}} \xrightarrow{\tilde{\mu}} 0$.

- Comment 3.96. Note $\widetilde{\boldsymbol{\mu}}$ depends on $S$, not just on a ratio set generating $S$.

For composition $T \circ S$, we need $S$ to be large. Example: Let $T=\sum_{j=0}^{\infty} x^{-j}$, $S=x^{-1}$. Then $S$ is small, not large. And $T \circ S=\sum_{j=0}^{\infty} x^{j}$ is not a valid transseries.
Remark 3.97. If $S \in \mathbb{T}_{N_{1}}$ and $T \in \mathbb{T}_{N_{2}}$, then $T \circ S \in \mathbb{T}_{N_{1}+N_{2}}$.
The height of $T \circ S$ is at most $N_{1}+N_{2}$. But sometimes the height of the composition is less than the sum of the heights. Note that if $n \leq N$, then $\mathfrak{G}_{n} \cup \mathfrak{G}_{N}^{\text {small }}=\mathfrak{G}_{n}^{\text {large }} \cup\{1\} \cup \mathfrak{G}_{N}^{\text {small }}$ is a convex subset of $\mathfrak{G}_{\boldsymbol{\bullet}}$.

Proposition 3.98. Let $N, n_{1}, n_{2} \in \mathbb{N}$ with $n_{1}+n_{2} \leq N$. Assume $T, S \in \mathbb{T} \bullet$, $\operatorname{supp} T \subset \mathfrak{G}_{n_{1}} \cup \mathfrak{G}_{N}^{\text {small }}, \operatorname{supp} S \subset \mathfrak{G}_{n_{2}} \cup \mathfrak{G}_{N}^{\text {small }}$. Then $\operatorname{supp}(T \circ S) \subset \mathfrak{G}_{n_{1}+n_{2}} \cup$ $\mathfrak{G}_{N}^{\text {small }}$.
Proof. We prove the result stated by induction on $n_{1}$. If $x^{b} \in \mathfrak{G}_{0}$, then write $T=c \mathfrak{g}+B$ where $\mathfrak{g}=\operatorname{mag} T, c \in \mathbb{R}, B \prec \mathfrak{g}$, so

$$
(c \mathfrak{g}+B)^{b}=c^{b} \mathfrak{g}^{b}\left(1+\frac{B}{c \mathfrak{g}}\right)^{b}=c^{b} \mathfrak{g}^{b} \sum_{j=0}^{\infty}\binom{b}{j}\left(\frac{B}{c \mathfrak{g}}\right)^{j} .
$$

This is supported by $\mathfrak{G}_{n_{2}} \cup \mathfrak{G}_{N}^{\text {small }}$. If the claim is true for all monomials in $\mathfrak{G}_{n_{1}}$ for some $n_{1}$, then summation shows it is still true for all transseries $T$ in $\mathbb{T}_{n_{1}}$. For the induction step: Suppose the claim is true for some $n_{1}<N-n_{2}$ and consider $\mathfrak{g} \in \mathfrak{G}_{n_{1}+1}$. So $\mathfrak{g}=x^{b} e^{L}, L$ purely large, $\operatorname{supp} L \subset \mathfrak{G}_{n_{1}}$. By the
induction hypothesis, $\operatorname{supp}(L \circ S) \subset \mathfrak{G}_{n_{1}+n_{2}} \cup \mathfrak{G}_{N}^{\text {small }}$. Write $L \circ S=U+c+V$, $U$ purely large, $c \in \mathbb{R}, V$ small. Then $\operatorname{supp} U \subseteq \mathfrak{G}_{n_{1}+n_{2}}^{\text {large }}$ and $\operatorname{supp} V \subset$ $\mathfrak{G}_{N}^{\text {small }}$. Now $e^{L \circ S}=e^{U} e^{c} e^{V}$, and $e^{U} \in \mathfrak{G}_{n_{1}+n_{2}+1}$, while $e^{V}=\sum_{j=0}^{\infty} V^{j} / j!$ is supported by $\{1\} \cup \mathfrak{G}_{N}^{\text {small }}$. This, together with $S^{b}$ done before, completes the induction.

## Continuity of Composition

Proposition 3.99. Let $S$ be large and positive. Let $\left(T_{i}\right)$ be a family of transseries with $T_{i} \rightarrow T$. Then $T_{i} \circ S \rightarrow T \circ S$.

Proof. Say $T_{i} \xrightarrow{\boldsymbol{\mu}, \mathbf{m}} T$. Let $\widetilde{\boldsymbol{\mu}}$ and $\widetilde{\mathbf{m}}$ be as in Proposition 3.95 so that $\mathfrak{g} \circ S \in$ $\mathbb{T}^{\tilde{\mu}, \tilde{m}}$ for all $\mathfrak{g} \in \mathfrak{J}^{\mu, \mathbf{m}}$, and $\mathfrak{g} \circ S \xrightarrow{\tilde{\mu}} 0$. Let $\mathfrak{m} \in \mathfrak{G}_{\boldsymbol{\bullet}}$. There are finitely many $\mathfrak{g} \in \mathfrak{J}^{\boldsymbol{\mu}, \mathbf{m}}$ such that $\mathfrak{m} \in \operatorname{supp}(\mathfrak{g} \circ S)$. For each such $\mathfrak{g}$ there are finitely many $i$ such that $\mathfrak{g} \in \operatorname{supp}\left(T_{i}-T\right)$. So if $i$ is outside this finite union of finite sets, we have $\mathfrak{m} \notin \operatorname{supp}\left(\left(T_{i}-T\right) \circ S\right)$.

- Comment 3.100. Continuity in the other composand might not hold. For $j \in \mathbb{N}$, let $S_{j}=x^{-j} e^{x}$. Then $\left(S_{j}\right) \rightarrow 0$. But the family $\left(\exp \left(S_{j}\right)\right)$ is not supported by any grid, so $\exp \left(S_{j}\right)$ cannot converge to anything.

Tedious calculation should show that the usual derivative formulas hold: $\left(S^{b}\right)^{\prime}=b S^{b-1} S^{\prime},\left(e^{T}\right)^{\prime}=e^{T} T^{\prime},(T \circ S)^{\prime}=\left(T^{\prime} \circ S\right) \cdot S^{\prime}$, and so on.

## 3E With Logarithms.

Transseries with logs are obtained by formally composing the log-free transseries with log on the right.
Notation 3.101. If $m \in \mathbb{N}$, we write $\log _{m}$ to represent the $m$-fold composition of the natural logarithm with itself; $\log _{0}$ will have no effect; sometimes we may write $\log _{m}=\exp _{-m}$, especially when $m<0$.
Definition 3.102. Let $M \in \mathbb{N}$. A transseries with depth $M$ is a formal expression $Q=T \circ \log _{M}$, where $T \in \mathbb{T}_{\bullet}$.

We identify the set of transseries of depth $M$ as a subset of the set of transseries of depth $M+1$ by identifying $T \circ \log _{M}$ with $(T \circ \exp ) \circ \log _{M+1}$. Composition on the right with exp is defined in Definition 3.93. Using this idea, we define operations on transseries from the operations in $\mathbb{T}_{\bullet}$.
Definition 3.103. Let $Q_{j}=T_{j} \circ \log _{M}$, where $T_{j} \in \mathbb{T}_{\bullet}$. Define $Q_{1}+Q_{2}=$ $\left(T_{1}+T_{2}\right) \circ \log _{M} ; Q_{1} Q_{2}=\left(T_{1} T_{2}\right) \circ \log _{M} ; Q_{1}>Q_{2}$ iff $T_{1}>T_{2} ; Q_{1} \succ Q_{2}$ iff $T_{1} \succ T_{2} ; Q_{j} \rightarrow Q_{0}$ iff $T_{j} \rightarrow T_{0} ; \sum Q_{j}=\left(\sum T_{j}\right) \circ \log _{M} ; Q_{1}^{b}=\left(T_{1}^{b}\right) \circ \log _{M} ;$ $\exp \left(Q_{1}\right)=\left(\exp \left(T_{1}\right)\right) \circ \log _{M}$; and so on.

Definition 3.104. Transseries. Always assumed grid-based.

$$
\begin{aligned}
& \mathfrak{G}_{N M}=\left\{\mathfrak{g} \circ \log _{M}: \mathfrak{g} \in \mathfrak{G}_{N}\right\}, \\
& \mathbb{T}_{N M}=\left\{T \circ \log _{M}: T \in \mathbb{T}_{N}\right\}=\mathbb{R} \llbracket \mathfrak{G}_{N M} \rrbracket, \\
& \mathfrak{G}_{\bullet M}=\bigcup_{N \in \mathbb{N}} \mathfrak{G}_{N M}=\left\{\mathfrak{g} \circ \log _{M}: \mathfrak{g} \in \mathfrak{G} \bullet\right\}, \quad \mathfrak{G}_{N \bullet}=\bigcup_{M \in \mathbb{N}} \mathfrak{G}_{N M}, \\
& \left.\mathbb{T}_{\bullet M}=\bigcup_{N \in \mathbb{N}} \mathbb{T}_{N M}=\left\{T \circ \log _{M}: T \in \mathbb{T}_{\bullet}\right\}=\mathbb{R} \llbracket \mathfrak{G} \bullet M \rrbracket, \quad \mathbb{T}_{N \bullet}=\mathbb{R} \llbracket \mathfrak{G}_{N \bullet} \rrbracket\right], \\
& \mathfrak{G}_{\bullet \bullet}=\bigcup_{M \in \mathbb{N}} \mathfrak{G}_{\bullet M}=\bigcup_{N, M \in \mathbb{N}} \mathfrak{G}_{N M}, \\
& \mathbb{R} \llbracket x \rrbracket \rrbracket=\bigcup_{M \in \mathbb{N}} \mathbb{T}_{\bullet M}=\bigcup_{N, M \in \mathbb{N}} \mathbb{T}_{N M}=\mathbb{R} \llbracket \mathfrak{G} \bullet \bullet \rrbracket=\mathbb{T} \bullet \bullet .
\end{aligned}
$$

When $M<0$ we also write $\mathbb{T}_{\bullet} \cdot$. So for example $\mathbb{T}_{\bullet,-1}=\left\{T \circ \exp : T \in \mathbb{T}_{\bullet}\right\}$.
If $T=\sum c_{\mathfrak{g}} \mathfrak{g}$ we may write $T \circ \log _{M}$ as a series

$$
\left(\sum c_{\mathfrak{g}} \mathfrak{g}\right) \circ \log _{M}=\sum c_{\mathfrak{g}}\left(\mathfrak{g} \circ \log _{M}\right)
$$

Simplifications along these lines may be carried out: $\exp (\log x)=x ; e^{b \log x}=$ $x^{b}$; etc. As usual we sometimes think of $x$ as a variable and sometimes as the identity function. On monomials we can write

$$
\left(x^{b} e^{L}\right) \circ \log =(\log x)^{b} e^{L \circ \log }
$$

and continue recursively in the exponent
We say $Q \in \mathbb{R} \llbracket x \rrbracket$ has exact depth $M$ iff $Q=T \circ \log _{M}, T \in \mathbb{T} \bullet$ and $T$ cannot be written in the form $T=T_{1} \circ \exp$ for $T_{1} \in \mathbb{T}_{\bullet}$. This will also make sense for negative $M$.

## Terminology.

The group $\mathfrak{G}=\mathfrak{G} .$. is the group of transmonomials. The ordered field $\mathbb{T}=\mathbb{T}_{\bullet \bullet}=\mathbb{R} \llbracket x \rrbracket$ is the field of transseries. Van der Hoeven [37] calls $\mathbb{T}$ the transline.

- Comment 3.105. Although $x^{x}$ is not an "official" transmonomial, if we consider it to be an abbreviation for $e^{x \log x}$, then it may be considered to be a transmonomial according to our identifications:

$$
x^{x}=e^{x \log x}=\left(e^{e^{x} x}\right) \circ \log
$$

So $x^{x}$ has height 2 and depth 1 ; that is, $x^{x} \in \mathfrak{G}_{2,1}$.

- Comment 3.106. Just as we require finite exponential height, we also require finite logarithmic depth. So the following is not a grid-based transseries:

$$
x+\log x+\log \log x+\log \log \log x+\cdots
$$

But see for example [34] or [23] for a variant that allows this.

- Comment 3.107. If $\mathfrak{g} \in \mathfrak{G}_{\bullet \bullet}$, then $\mathfrak{g}=e^{L}$ for some purely large $L \in \mathbb{T}_{\bullet \bullet}$. Because of logarithms, there is no need for an extra $x^{b}$ factor.


## Composition.

Definition 3.108. Logarithm. If $T \in \mathbb{T}_{\bullet}, T>0$, write $T=a x^{b} e^{L}(1+S)$ as in 3.29. Define $\log T=\log a+b \log x+L+\log (1+S)$. Now $\log a, a>0$, is computed in $\mathbb{R}$. And $\log (1+S)$ is computed as a Taylor series. The term $b \log x$ gives this depth 1 ; if $b=0$ then we remain log-free.

For general $Q \in \mathbb{R} \mathbb{\llbracket} x \mathbb{\square}$ : if $Q=T \circ \log _{M}$, then $\log (Q)=\log (T) \circ \log _{M}$, which could have depth $M+1$.

Alternatively (from Comment 3.107): for $Q \in \mathbb{T} \bullet \bullet$, write $Q=a e^{L}(1+S)$ and then $\log Q=L+\log a+\log (1+S)$.

- Comment 3.109. If $T \not \not 1$, then $\log T \succ 1$. This is because in the definition, $\log (1+S)=S+\cdots$ is small.

Definition 3.110. Composition. Let $Q_{1}, Q_{2} \in \mathbb{R} \llbracket x \rrbracket$ with $Q_{2}$ large and positive. Define $Q_{1} \circ Q_{2}$ as follows: Write $Q_{1}=T_{1} \circ \log _{M_{1}}$ and $Q_{2}=T_{2} \circ \log _{M_{2}}$, with $T_{1}, T_{2} \in \mathbb{T}_{\bullet}$. Applying $3.108 M_{1}$ times, we can write $\log _{M_{1}}\left(T_{2}\right)=S \circ \log _{M_{1}}$ with $S \in \mathbb{T}_{\bullet}$. Then define:

$$
Q_{1} \circ Q_{2}=T_{1} \circ \log _{M_{1}} \circ T_{2} \circ \log _{M_{2}}=T_{1} \circ S \circ \log _{M_{1}+M_{2}},
$$

and compute $T_{1} \circ S$ as in 3.93.
The depth of $Q_{1} \circ Q_{2}$ is at most $M_{1}+M_{2}$. But sometimes the depth of the composition is less than the sum of the depths.

Proposition 3.111. Let $N, M \in \mathbb{N}, N \geq M$.
(a) if $\operatorname{supp} A \subset \mathfrak{G}_{N}, A \preccurlyeq 1$, then $\log \left(\exp _{M}+A\right)=\exp _{M-1}+B$ with $\operatorname{supp} B \subset \mathfrak{G}_{N}, B \prec 1$.
(b) if $\operatorname{supp} A \subset \mathfrak{G}_{N}, A \preccurlyeq 1$, then $\log _{M}\left(\exp _{M}+A\right)=x+B$ with $\operatorname{supp} B \subset$ $\mathfrak{G}_{N}, B \prec 1$.
(c) if $\operatorname{supp} Q \subset \mathfrak{G}_{N M}$, $\operatorname{supp} A \subset \mathfrak{G}_{N M}, A \preccurlyeq 1$, then $Q \circ(x+A)=B$, with $\operatorname{supp} B \subset \mathfrak{G}_{N M}$.

Proof. (a) $\log \left(\exp _{M}+A\right)=\log \left(\exp _{M} \cdot\left(1+A / \exp _{M}\right)\right)$. But $\operatorname{supp}\left(A / \exp _{M}\right) \subset$ $\mathfrak{G}_{N}^{\text {small }}$ so $\log \left(\exp _{M}+A\right)$ is $\exp _{M-1}$ plus a series in positive powers of $A / \exp _{M}$.
(b) Apply (a) with induction.
(c) We have $Q=Q_{1} \circ \log _{M}, \operatorname{supp}\left(Q_{1}\right) \subset \mathfrak{G}_{N} ; A=A_{1} \circ \log _{M}, \operatorname{supp}\left(A_{1}\right) \subset$ $\mathfrak{G}_{N}$. Then

$$
\begin{aligned}
Q \circ(x+A) & =Q_{1} \circ \log _{M} \circ\left(x+\left(A_{1} \circ \log _{M}\right)\right) \\
& =Q_{1} \circ \log _{M} \circ\left(\exp _{M}+A_{1}\right) \circ \log _{M} .
\end{aligned}
$$

By (b), $\log _{M} \circ\left(\exp _{M}+A_{1}\right)=x+B_{1}$ with $\operatorname{supp}\left(B_{1}\right) \subset \mathfrak{G}_{N}$. So

$$
B:=Q \circ(x+A)=Q_{1} \circ\left(x+B_{1}\right) \circ \log _{M} .
$$

By Proposition 3.98, $\operatorname{supp}\left(Q_{1} \circ\left(x+B_{1}\right)\right) \subset \mathfrak{G}_{N}$. Therefore $\operatorname{supp} B \subset \mathfrak{G}_{N M}$.
The set of large positive transseries from $\mathbb{R} \llbracket x \rrbracket$ is closed under composition. In fact, it is a group [37, p. 111].

## Derivative.

Definition 3.112. Differentiation is done as expected from the usual rules.

$$
(T \circ \log )^{\prime}=\left(T^{\prime} \circ \log \right) \cdot x^{-1}=\left(T^{\prime} e^{-x}\right) \circ \log
$$

So $\partial$ maps $\mathbb{T}_{\bullet} \cdot M$ into itself.

- Comment 3.113. ... but perhaps $\partial$ does not map $\mathbb{T}_{N M}$ into itself. Example: $Q=(\log x)^{2}=\left(x^{2}\right) \circ \log \in \mathbb{T}_{0,1}$, and $Q^{\prime}=2(\log x) / x=\left(2 x e^{-x}\right) \circ \log \in \mathbb{T}_{1,1}$ but $Q^{\prime} \notin \mathbb{T}_{0,1}$. See Remark 3.116.

We now have an antiderivative for $x^{-1}$.

$$
(\log x)^{\prime}=(x \circ \log )^{\prime}=\left(1 \cdot e^{-x}\right) \circ \log =\left(x^{-1}\right) \circ \exp \circ \log =x^{-1}
$$

We will see below (Proposition 4.29) that, in fact, every transseries has an antiderivative.

Here are some simple properties of the derivative.
Proposition 3.114. Let $S, T \in \mathbb{R} \llbracket x \rrbracket$, $\mathfrak{n}, \mathfrak{m} \in \mathfrak{G}$... (a) If $\mathfrak{m} \succ \mathfrak{n}$, and $\mathfrak{m} \neq 1$, then $\mathfrak{m}^{\prime} \succ \mathfrak{n}^{\prime}$. (b) If $\operatorname{mag} T \neq 1$, then $T^{\prime} \asymp(\operatorname{mag} T)^{\prime}$ and $T^{\prime} \sim(\operatorname{dom} T)^{\prime}$. (c) If $\operatorname{mag} T \neq 1$ and $T \succ S$, then $T^{\prime} \succ S^{\prime}$. (d) If $T \prec 1$, then $T^{\prime} \prec 1$ 。 (e) If $T \succ 1$ and $T>0$, then $T^{\prime}>0$. (f) If $T \succ 1$ and $T<0$, then $T^{\prime}<0$. (g) If $T \succ 1$, then $T^{2} \succ T^{\prime}$. (h) If $T \succ \log x$, then $\left(T^{\prime}\right)^{2} \succ T^{\prime \prime}$. (i) If $T \succ \log x$, then $x T^{\prime} \succ 1$. ( j ) If $T^{\prime}=0$, then $T$ is a constant.

Proof. (a)(b)(c) Starting with Proposition 3.81(b)(c)(d), compose with log repeatedly.
(d) $T \prec 1 \quad \Longrightarrow \quad T \prec x \quad \Longrightarrow \quad T^{\prime} \prec 1$.
(e)(f) Starting with Proposition 3.82 (ii)(iii), compose with log repeatedly.
(g) $T \succ 1 \Longrightarrow T \succ \frac{1}{x} \Longrightarrow \frac{1}{T} \prec x \Longrightarrow\left(\frac{1}{T}\right)^{\prime} \prec 1 \Longrightarrow \frac{T^{\prime}}{T^{2}} \prec 1$.
(h) since $T \succ \log x$, we have $T^{\prime} \succ x^{-1}$, then proceed as in (g).
(i) $T \succ \log x \quad \Longrightarrow \quad T^{\prime} \succ x^{-1}$.
(j) Starting with Proposition 3.85, compose with log repeatedly.

- Comment 3.115. Note $T=\log \log x$ is a counterexample to: If $T \succ 1$ then $x T^{\prime} \succ 1$. And to: If $T \succ 1$, then $\left(T^{\prime}\right)^{2} \succ T^{\prime \prime}$.

Remark 3.116. Of course $\mathfrak{G}_{N M}$ is a group, so $\mathbb{T}_{N M}$ is a field. The derivative of $\log _{M}$ is

$$
\left(\log _{M} x\right)^{\prime}=\left(\prod_{m=0}^{M-1} \log _{m} x\right)^{-1}
$$

If $N \geq M$ then this belongs to $\mathfrak{G}_{N M}$, so in that case $\mathbb{T}_{N M}$ is a differential field (with constants $\mathbb{R}$ ).

Proposition 3.117. If $T \in \mathbb{T}_{N, M}$ has exact height $N \geq 1$ (that is, $T \notin$ $\left.\mathbb{T}_{N-1, \bullet}\right)$ and $T^{\prime} \succ 1$, then $T^{\prime} \in \mathbb{T}_{N, M}$ and also has exact height $N$.

Proof. There is $T_{1} \in \mathbb{T}_{N}$ with $T(x)=T_{1}\left(\log _{M} x\right)$. Then

$$
T^{\prime}(x)=\frac{T_{1}^{\prime}\left(\log _{M} x\right)}{x \log x \cdots \log _{M-1} x}
$$

so $T_{1}^{\prime}\left(\log _{M} x\right) \succ x$ and $T_{1}(x) \succ \exp _{M} x$. So $N \geq M$ and $x \log x \cdots \log _{M-1} x \in$ $\mathbb{T}_{N-1, M}$ so $T^{\prime} \in \mathbb{T}_{N, M}$. Since $T_{1}^{\prime}$ has exact height $N$ and $x \log x \cdots \log _{M-1} x$ has height $\leq N-1$, it follows that $T^{\prime}$ has exact height $N$.

## Valuation.

The map "mag" from $\mathbb{R} \llbracket \mathfrak{G} \rrbracket \backslash\{0\}$ to $\mathfrak{G}$ is a (nonarchimedean) valuation. This means:
(i) $\operatorname{mag}(S T)=\operatorname{mag}(S) \operatorname{mag}(T)$;
(ii) $\operatorname{mag}(S+T) \preccurlyeq \max \{\operatorname{mag}(S), \operatorname{mag}(T)\}$ with equality if $\operatorname{mag}(S) \neq \operatorname{mag}(T)$.

The ordered group $\mathfrak{G}$ is the valuation group.

- Comment 3.118. The valuation group is written multiplicatively here, but in many parts of mathematics it is more common to write it additively, and with the order reversed. In the transseries case, $\mathfrak{G}=\mathfrak{G}$ ••, we could follow this "additive" convention by saying: the valuation group is the set of purely large transseries, with operation + and order $<$. The valuation $\nu$ is then related to the magnitude by: mag $T=e^{-L} \Longleftrightarrow \nu(T)=L$. We could then still call $\mathfrak{G}$ the monomial group. But for a general ordered abelian monomial group $\mathfrak{G}$ (without log and exp) the valuation group would have to consist of "formal logarithms" of the monomials; introducing them may seem artificial. -

The map "mag" is an ordered valuation. This means that it also satisfies:
(iii) if $\operatorname{mag}(T) \succ 1$ then $|T|>1$. [The absolute value $|T|$ is defined as usual.]

The map "mag" is a differential valuation. This means that it also satisfies:
(iv) if $\operatorname{mag}(T) \neq 1, \operatorname{mag}(S) \neq 1$, then $\operatorname{mag}(T) \preccurlyeq \operatorname{mag}(S)$ if and only if $\operatorname{mag}\left(T^{\prime}\right) \preccurlyeq \operatorname{mag}\left(S^{\prime}\right) ;$
(v) if $\operatorname{mag}(T) \prec \operatorname{mag}(S) \neq 1$, then $\operatorname{mag}\left(T^{\prime}\right) \prec \operatorname{mag}\left(S^{\prime}\right)$.

For more on valuations, see [43] and [50].

## 4 Example Computations.

I will show here some computations. They can be done by hand with patience, but modern computer algebra systems will handle them easily. Read these or-better yet-try doing some computations of your own. I think that your own experience with it will convince you better than anything else that this system is truly elementary, but very powerful.

## A Polynomial Equation.

Problem 4.1. Solve the fifth-degree polynomial equation

$$
P(Y):=Y^{5}+e^{x} Y^{2}-x Y-9=0
$$

for $Y$.
We can think of this problem in various ways. If $x$ is a real number, then we want to solve for a real number $Y$. (When $x=0$, the Galois group is $S_{5}$, so we will not be solving this by radicals!) Or: think of $x$ and $e^{x}$ as functions, then the solution $Y$ is to be a function as well. Or: think of $x$ and $e^{x}$ as
transseries, then the solution $Y$ is to be a transseries. From some points of view, this last one is the easiest of the three. That is what we will do now.

In fact, many polynomial equations with transseries coefficients have transseries solutions. Of course for solutions in $\mathbb{R} \llbracket x \mathbb{\square}$ there are certain restrictions, since some polynomials (such as $Y^{2}+1$ ) might have no zeros because they are always positive. But if $P(Y) \in \mathbb{R} \mathbb{T} x \mathbb{\rrbracket}[Y]$ and there are transseries $A, B$ with $P(A)>0, P(B)<0$, there there is a transseries $T$ between $A$ and $B$ with $P(T)=0$. Van der Hoeven [37, Chap. 9] has this is even for differential polynomials.

The transseries coefficients of our equation $Y^{5}+e^{x} Y^{2}-x Y-9=0$ belong to the set $\mathbb{T}_{1}=\mathbb{T}_{1,0}$ of height 1 depth 0 transseries. Our solutions will also be in $\mathbb{T}_{1}$. Note that $Y=0$ is not a solution. So any solution $Y$ has a dominance $\operatorname{dom} Y=a x^{b} e^{L}$, where $a \neq 0$ and $b$ are real, and $L \in \mathbb{T}_{0}$ is purely large. So the dominances for the terms are:

$$
\begin{aligned}
\operatorname{dom}\left(Y^{5}\right) & =a^{5} x^{5 b} e^{5 L} \\
\operatorname{dom}\left(e^{x} Y^{2}\right) & =a^{2} x^{2 b} e^{2 L+x} \\
\operatorname{dom}(-x Y) & =-a x^{b+1} e^{L} \\
\operatorname{dom}(-9) & =-9
\end{aligned}
$$

Now we can compare these four terms. If $L>x / 3$, then $Y^{5}$ is far larger than any other term, so $P(Y) \asymp Y^{5}$ and therefore $P(Y) \neq 0$. If $-x / 2<L<x / 3$, then $P(Y) \asymp e^{x} Y^{2}$. If $L<-x / 2$, then $P(T) \asymp-9$ (all other terms are $\prec 1$ ). So the only possibilities for $L$ are $x / 3$ and $-x / 2$. [If you know the "Newton polygon" method, you may recognize what we just did.]

We consider first $L=x / 3$. If $b>0$ then $P(Y) \asymp Y^{5}$; if $b<0$ then $P(Y) \asymp e^{x} Y^{2}$; so $b=0$. Then we must have $\operatorname{dom}\left(Y^{5}\right)+\operatorname{dom}\left(e^{-x} Y^{2}\right)=0$, since otherwise the sum $P(Y) \asymp Y^{5}+e^{-x} Y^{2}$. So $a^{5}+a^{2}=0$. Since $a \neq 0$ and $a \in \mathbb{R}$, we have $a=-1$. [To consider also complex zeros of $P$, we would try to use complex-valued transseries, and then the other two cube roots of -1 would also need to be considered here.] Thus $Y=-e^{x / 3}(1+S)$, where $S \prec 1$. Then $P\left(-e^{x / 3}(1+S)\right)=$

$$
\begin{aligned}
\left(-e^{x / 3}(1+S)\right)^{5}+ & e^{x}\left(-e^{x / 3}(1+S)\right)^{2}-x\left(-e^{x / 3}(1+S)\right)-9 \\
= & -3 e^{5 x / 3} S-9 e^{5 x / 3} S^{2}-10 e^{5 x / 3} S^{3} \\
& \quad-5 e^{5 x / 3} S^{4}-e^{5 x / 3} S^{5}+x e^{x / 3} S+x e^{x / 3}-9
\end{aligned}
$$

Since $S$ is small, among the terms involving $S$ the dominant one is $-3 e^{5 x / 3} S$. Solve $P\left(-e^{x / 3}(1+S)\right)=0$ for that term, and write the equation as $S=\Phi(S)$,
where

$$
\Phi(S):=-3 S^{2}-\frac{10}{3} S^{3}-\frac{5}{3} S^{4}-\frac{1}{3} S^{5}+\frac{1}{3} x e^{-4 x / 3} S+\frac{1}{3} x e^{-4 x / 3}-3 e^{-5 x / 3}
$$

Start with any $S_{0}$ and iterate $S_{1}=\Phi\left(S_{0}\right), S_{2}=\Phi\left(S_{1}\right)$, etc. For example,

$$
\begin{aligned}
S_{0} & =0 \\
S_{1} & =\frac{1}{3} x e^{-4 x / 3}-3 e^{-5 x / 3} \\
S_{2} & =\frac{1}{3} x e^{-4 x / 3}-3 e^{-5 x / 3}-\frac{2}{9} x^{2} e^{-8 x / 3}+5 x e^{-3 x}-27 e^{-10 x / 3}-\frac{10}{81} x^{3} e^{-4 x} \ldots \\
S_{3} & =\frac{1}{3} x e^{-4 x / 3}-3 e^{-5 x / 3}-\frac{2}{9} x^{2} e^{-8 x / 3}+5 x e^{-3 x}-27 e^{-10 x / 3}+\frac{20}{81} x^{3} e^{-4 x} \ldots
\end{aligned}
$$

Each step produces more terms that subsequently remain unchanged. Thus we get a solution for $P(Y)=0$ in the form

$$
\begin{aligned}
Y= & -e^{x / 3}-\frac{1}{3} x e^{-x}+3 e^{-4 x / 3}+\frac{2}{9} x^{2} e^{-7 x / 3}-5 x e^{-8 x / 3}+27 e^{-3 x} \\
& -\frac{20}{81} x^{3} e^{-11 x / 3}+9 x^{2} e^{-4 x}-105 x e^{-13 x / 3}+396 e^{-14 x / 3}+\frac{1}{3} x^{4} e^{-5 x} \\
& -\frac{455}{27} x^{3} e^{-16 x / 3}+308 x^{2} e^{-17 x / 3}-2430 x e^{-6 x}-\frac{364}{729} x^{5} e^{-19 x / 3} \\
& +7020 e^{-19 x / 3}+\frac{2618}{81} x^{4} e^{-20 x / 3}-810 x^{3} e^{-7 x}+\mathrm{o}\left(e^{-7 x}\right)
\end{aligned}
$$

The "little o" on the end represents, as usual, a remainder that is $\prec e^{-7 x}$.
Now consider the other possibility, $L=-x / 2$. Using the same reasoning as before, we get $b=0$ and $a^{2}-9=0$, so there are two possibilities, $a= \pm 3$. With the same steps as before, we end up with two more solutions,

$$
\begin{aligned}
Y= & \pm 3 e^{-x / 2}+\frac{1}{2} x e^{-x} \pm \frac{1}{24} x^{2} e^{-3 x / 2} \mp \frac{1}{3456} x^{4} e^{-5 x / 2}-\frac{81}{2} e^{-3 x} \\
& \pm \frac{1}{248832} x^{6} e^{-7 x / 2} \mp \frac{135}{4} x e^{-7 x / 2}-\frac{27}{2} x^{2} e^{-4 x} \mp \frac{5}{71663616} x^{8} e^{-9 x / 2} \\
& \mp \frac{105}{32} x^{3} e^{-9 x / 2}-\frac{1}{2} x^{4} e^{-5 x} \pm \frac{7}{5159780352} x^{10} e^{-11 x / 2} \mp \frac{21}{512} x^{5} e^{-11 x / 2} \\
& \pm \frac{19683}{8} e^{-11 x / 2}+3645 x e^{-6 x} \mp \frac{7}{247669456896} x^{12} e^{-13 x / 2} \\
& \pm \frac{11}{36864} x^{7} e^{-13 x / 2} \pm \frac{168399}{64} x^{2} e^{-13 x / 2}+1215 x^{3} e^{-7 x}+\mathrm{o}\left(e^{-7 x}\right)
\end{aligned}
$$

It turns out that these three transseries solutions converge for large enough $x$. As a check, take $x=10$ in $P$. Maple says the zeros are

$$
\begin{aligned}
& -28.0317713673296286443879064009, \\
& \quad-0.0199881159048462608264265543923, \\
& \quad 0.0204421151948799622524221088662, \\
& 14.0156586840197974714809554232+24.27610347738050183477184 i, \\
& 14.0156586840197974714809554232-24.27610347738050183477184 i .
\end{aligned}
$$

Plugging $x=10$ in the three series shown above (up to order $e^{-7 x}$ ), I get

$$
\begin{array}{r}
-28.0317713673296286443879064149 \\
-0.0199881159048462608264265439647 \\
0.0204421151948799622524220981049
\end{array}
$$

## A Derivative and a Borel Summation.

Consider the Euler series

$$
S=\sum_{k=0}^{\infty} \frac{k!e^{x}}{x^{k+1}}
$$

Differentiate term-by-term to get a telescoping sum leaving only $S^{\prime}=e^{x} / x$. This means any summation method that commutes with summation of series and differentiation should yield the exponential integral function for $S$. In fact, this is a case that can be done by classical Borel summation (Example 2.10).

## A Compositional Inverse

Problem 4.2. Compute the compositional inverse of $x e^{x}$
This inverse is known as the Lambert $W$ function. There is a standard construction for all compositional inverses, but we will proceed here directly. First we need to know the dominant term. This is done by "reducing to height zero" as follows. If $x=W e^{W}$, then $\log x=W+\log W$, so $W=\log x-\log W$ with $\log W \prec W$, and thus $W \sim \log x$.

So assume our inverse is $\log x+Q$, with $Q \prec \log x$. Then $x=(\log x+$ $Q) e^{\log x+Q}$ so $x=(\log x+Q) x e^{Q}$ so $e^{-Q}=\log x+Q$. Now we should solve for one $Q$ in terms of the other(s), and use this to iterate. If we take $Q=$ $e^{-Q}-\log x$ and iterate $\Phi(Q)=e^{-Q}-\log x$, it doesn't work: starting with $Q_{0}=0$, we get $Q_{1}=1-\log x$, then $Q_{2}=x / e-\log x$, which is not converging.

So we will solve for the other one: $Q=-\log (\log x+Q)$. Write $\Phi(Q)=$ $-\log (\log x+Q)$ and iterate. Since we assume $Q \prec \log x$, the term $Q / \log x$ is small. So write

$$
\begin{aligned}
\Phi(Q) & =-\log (\log x+Q)=-\log ((\log x)(1+Q / \log x)) \\
& =-\log \log x+\sum_{k=1}^{\infty} \frac{(-1)^{k}}{k}\left(\frac{Q}{\log x}\right)^{k}
\end{aligned}
$$

We will start with ratios $\mu_{1}=1 / \log \log x, \mu_{2}=1 / \log x$. So

$$
\Phi(Q)=-\mu_{1}^{-1}+\sum_{k=1}^{\infty} \frac{(-1)^{k}}{k}\left(Q \mu_{2}\right)^{k}
$$

Start with $Q_{0}=0$. Then $Q_{1}$ begins $-\mu_{1}^{-1}$, so for the series in $\Phi\left(Q_{1}\right)$ we need $\mu_{1}^{-1} \mu_{2} \prec 1$. Of course $\mu_{1}^{-1} \mu_{2}=\log \log x / \log x$ actually is small, but not $\left\{\mu_{1}, \mu_{2}\right\}$-small. So we add another ratio, $\mu_{3}=\log \log x / \log x=\mu_{1}^{-1} \mu_{2}$. Now computing with $\mu_{2}$ and $\mu_{3}$, iteration of

$$
\Phi(Q)=-\mu_{2}^{-1} \mu_{3}+\sum_{k=1}^{\infty} \frac{(-1)^{k}}{k}\left(Q \mu_{2}\right)^{k}
$$

is just a matter of routine:

$$
\begin{aligned}
Q_{0} & =0 \\
Q_{1} & =-\mu_{2}^{-1} \mu_{3} \\
Q_{2} & =-\mu_{2}^{-1} \mu_{3}+\mu_{3}+\frac{1}{2} \mu_{3}^{2}+\frac{1}{3} \mu_{3}^{3}+\frac{1}{4} \mu_{3}^{4}+\frac{1}{5} \mu_{3}^{5}+\ldots \\
Q_{3} & =-\mu_{2}^{-1} \mu_{3}+\left(1-\mu_{2}\right) \mu_{3}+\left(\frac{1}{2}-\frac{3}{2} \mu_{2}+\frac{1}{2} \mu_{2}^{2}\right) \mu_{3}^{2}+\ldots
\end{aligned}
$$

When we continue this, we get more and more terms which remain the same from one step to the next. I did this with Maple, keeping terms with total degree at most 6 in $\mu_{2}, \mu_{3}$. When it stops changing, I know I have the first terms of the answer. Substituting in the values for $\mu_{2}$ and $\mu_{3}$, writing $\mathfrak{l}_{1}=\log x$ and $\mathfrak{l}_{2}=\log \log x$, we get:

$$
\begin{aligned}
W(x)= & \mathfrak{l}_{1}-\mathfrak{l}_{2}+\frac{\mathfrak{l}_{2}}{\mathfrak{l}_{1}}+\frac{1}{2} \frac{\mathfrak{l}_{2}^{2}}{\mathfrak{l}_{1}^{2}}-\frac{\mathfrak{l}_{2}}{\mathfrak{l}_{1}^{2}}+\frac{1}{3} \frac{\mathfrak{l}_{2}^{3}}{\mathfrak{l}_{1}^{3}}-\frac{3}{2} \frac{\mathfrak{l}_{2}^{2}}{\mathfrak{l}_{1}^{3}}+\frac{\mathfrak{l}_{2}}{\mathfrak{l}_{1}^{3}}+\frac{1}{4} \frac{\mathfrak{l}_{2}^{4}}{\mathfrak{l}_{1}^{4}}-\frac{11}{6} \frac{\mathfrak{l}_{2}^{3}}{\mathfrak{l}_{1}^{4}}+\frac{3 \mathfrak{l}_{2}^{2}}{\mathfrak{l}_{1}^{4}} \\
& -\frac{\mathfrak{l}_{2}}{\mathfrak{l}_{1}^{4}}+\frac{1}{5} \frac{\mathfrak{l}_{2}^{5}}{5}-\frac{25}{12} \frac{\mathfrak{l}_{2}^{4}}{\mathfrak{l}_{1}^{5}}+\frac{35}{6} \frac{\mathfrak{l}_{2}^{3}}{\mathfrak{l}_{1}^{5}}-\frac{5 \mathfrak{l}_{2}^{2}}{\mathfrak{l}_{1}^{5}}+\frac{\mathfrak{l}_{2}}{\mathfrak{l}_{1}^{5}}+\frac{1}{6} \frac{\mathfrak{l}_{2}^{6}}{\mathfrak{l}_{1}^{6}}-\frac{137}{60} \frac{\mathfrak{l}_{2}^{5}}{\mathfrak{l}_{1}^{6}}+\frac{75}{8} \frac{\mathfrak{l}_{2}^{4}}{\mathfrak{l}_{1}^{6}} \\
& -\frac{85}{6} \frac{\mathfrak{l}_{2}^{3}}{\mathfrak{l}_{1}^{6}}+\frac{15}{2} \frac{\mathfrak{l}_{2}^{2}}{\mathfrak{l}_{1}^{6}}-\frac{\mathfrak{l}_{2}}{\mathfrak{l}_{1}^{6}}+\cdots .
\end{aligned}
$$

## Contractive Mappings.

There is a general principle that explains why the sort of iterations that we have seen will work. It is a sort of "fixed-point" theorem for an appropriate type of "contraction" mappings. Here is an explanation.

First consider a domination relation for sets of multi-indices.
Definition 4.3. Let $\mathbf{E}, \mathbf{F}$ be subsets of $\mathbb{Z}^{n}$. We say $\mathbf{E}$ dominates $\mathbf{F}$ iff for every $\mathbf{k} \in \mathbf{F}$, there is $\mathbf{p} \in \mathbf{E}$ with $\mathbf{p}<\mathbf{k}$.

This may seem backward. But correspondingly in the realm of transmonomials, we will say larger monomials dominate smaller ones.

It's transitive: If $\mathbf{E}_{1}$ dominates $\mathbf{E}_{2}$ and $\mathbf{E}_{2}$ dominates $\mathbf{E}_{3}$, then $\mathbf{E}_{1}$ dominates $\mathbf{E}_{3}$. Every $\mathbf{E}$ dominates $\varnothing$.

Recall (Proposition 3.7) that Min $\mathbf{E}$ denotes the set of minimal elements of $\mathbf{E}$. And $\operatorname{Min} \mathbf{E}$ is finite if $\mathbf{E} \subseteq \mathbf{J}_{\mathbf{m}}$ for some $\mathbf{m}$.

Proposition 4.4. Let $\mathbf{E}, \mathbf{F}$ be subsets of $\mathbf{J}_{\mathbf{m}}$. Then $\mathbf{E}$ dominates $\mathbf{F}$ if and only if $\operatorname{Min} \mathbf{E}$ dominates $\operatorname{Min} \mathbf{F}$.

Proof. Assume $\mathbf{E}$ dominates $\mathbf{F}$. Let $\mathbf{k} \in \operatorname{Min} \mathbf{F}$. Then $\mathbf{k} \in \mathbf{F}$, so there is $\mathbf{k}_{1} \in \mathbf{E}$ with $\mathbf{k}_{1}<\mathbf{k}$. Then there is $\mathbf{k}_{0} \in \operatorname{Min} \mathbf{E}$ with $\mathbf{k}_{0} \leq \mathbf{k}_{1}$. So $\mathbf{k}_{0}<\mathbf{k}$.

Conversely, assume $\operatorname{Min} \mathbf{E}$ dominates $\operatorname{Min} \mathbf{F}$. Let $\mathbf{k} \in \mathbf{F}$. Then there is $\mathbf{k}_{1} \in \operatorname{Min} \mathbf{F}$ with $\mathbf{k}_{1} \leq \mathbf{k}$. So there is $\mathbf{k}_{0} \in \operatorname{Min} \mathbf{E}$ with $\mathbf{k}_{0}<\mathbf{k}_{1}$. Thus $\mathbf{k}_{0} \in \mathbf{E}$ and $\mathbf{k}_{0}<\mathbf{k}$.

Proposition 4.5. Let $\mathbf{E}, \mathbf{F} \subseteq \mathbf{J}_{\mathbf{m}}$. If $\mathbf{E}$ dominates $\mathbf{F}$, then $\operatorname{Min} \mathbf{E}$ and $\operatorname{Min} \mathbf{F}$ are disjoint.

Proof. Assume $\mathbf{E}$ dominates $\mathbf{F}$. If $\mathbf{k} \in \operatorname{Min} \mathbf{F}$, then $\mathbf{k} \in \mathbf{F}$, so there is $\mathbf{k}_{1} \in \mathbf{E}$ with $\mathbf{k}_{1}<\mathbf{k}$. So even if $\mathbf{k} \in \mathbf{E}$, it is not minimal.

Proposition 4.6. Let $\mathbf{E}_{j} \subseteq \mathbf{J}_{\mathbf{m}}, j \in \mathbb{N}$, be an infinite sequence such that $\mathbf{E}_{j}$ dominates $\mathbf{E}_{j+1}$ for all $j$. Then the sequence $\left(\mathbf{E}_{j}\right)$ is point-finite; $\mathbf{E}_{j} \rightarrow \varnothing$.

Proof. Let $\mathbf{p} \in \mathbf{J}_{\mathbf{m}}$. Then $\mathbf{F}=\left\{\mathbf{k} \in \mathbf{J}_{\mathbf{m}}: \mathbf{k} \leq \mathbf{p}\right\}$ is finite. But the sets $\mathbf{F} \cap \operatorname{Min} \mathbf{E}_{j}$ are disjoint (by Proposition 4.5), so all but finitely many of them are empty. For every $j$ with $\mathbf{p} \in \mathbf{E}_{j}$, the set $\mathbf{F} \cap \operatorname{Min} \mathbf{E}_{j}$ is nonempty. Therefore, $\mathbf{p} \in \mathbf{E}_{j}$ for only finitely many $j$.

Proposition 4.7. Let $\mathbf{E}_{i} \subseteq \mathbf{J}_{\mathbf{m}}$ be a point-finite family. Assume $\mathbf{E}_{i}$ dominates $\mathbf{F}_{i}$ for all $i$. Then $\left(\mathbf{F}_{i}\right)$ is also point-finite.

Proof. Let $\mathbf{p} \in \mathbf{J}_{\mathbf{m}}$. Then $\mathbf{F}=\left\{\mathbf{k} \in \mathbf{J}_{\mathbf{m}}: \mathbf{k}<\mathbf{p}\right\}$ is finite. But the collection of sets $\mathbf{F} \cap \operatorname{Min} \mathbf{E}_{j}$ is point-finite, so again all but finitely many of them are nonempty. For every $j$ with $\mathbf{p} \in \mathbf{F}_{j}$, the set $\mathbf{F} \cap \operatorname{Min} \mathbf{E}_{j}$ is nonempty. Therefore, $\mathbf{p} \in \mathbf{F}_{j}$ for only finitely many $j$.

Next consider the corresponding notion for a grid-based field $\mathbb{T}^{\mu}$ of transseries.
Definition 4.8. For $\mathfrak{m}, \mathfrak{n} \in \mathfrak{J}^{\boldsymbol{\mu}}$, we write $\mathfrak{m} \prec^{\boldsymbol{\mu}} \mathfrak{n}$ and we say $\mathfrak{n} \boldsymbol{\mu}$-dominates $\mathfrak{m}$ iff $\mathfrak{m} / \mathfrak{n} \prec^{\boldsymbol{\mu}} 1$ (that is, $\mathfrak{m} / \mathfrak{n}=\boldsymbol{\mu}^{\mathbf{k}}$ for some $\mathbf{k}>\mathbf{0}$ ).

The following are easy. (They follow from Propositions 3.5-3.7 using Proposition 3.42). The grid $\mathfrak{J}^{\mu, \mathbf{m}}$ is well-partially-ordered for (the converse of) $\succ^{\mu}$.

Proposition 4.9. If $\mathfrak{A} \subseteq \mathfrak{J}^{\boldsymbol{\mu}, \mathbf{m}}, \mathfrak{A} \neq \varnothing$, then there is a $\boldsymbol{\mu}$-maximal element: $\mathfrak{m} \in \mathfrak{A}$ and $\mathfrak{g} \succ^{\mu} \mathfrak{m}$ for no $\mathfrak{g} \in \mathfrak{A}$.

Proposition 4.10. Let $\mathfrak{A} \subseteq \mathfrak{J}^{\mu, \mathbf{m}}$ be infinite. Then there is a sequence $\mathfrak{g}_{j} \in \mathfrak{A}$, $j \in \mathbb{N}$, with $\mathfrak{g}_{0} \succ^{\mu} \mathfrak{g}_{1} \succ^{\mu} \mathfrak{g}_{2} \succ^{\mu} \ldots$.
Proposition 4.11. Let $\mathfrak{A} \subseteq \mathfrak{J}^{\boldsymbol{\mu , m}}$. Then the set $\operatorname{Max}^{\boldsymbol{\mu}} \mathfrak{A}$ of $\boldsymbol{\mu}$-maximal elements of $\mathfrak{A}$ is finite. For every $\mathfrak{g} \in \mathfrak{A}$ there is $\mathfrak{m} \in \operatorname{Max}^{\mu} \mathfrak{A}$ with $\mathfrak{g} \preccurlyeq^{\mu} \mathfrak{m}$.
Definition 4.12. Let $\mathfrak{A}, \mathfrak{B} \subseteq \mathfrak{G}$. We say $\mathfrak{A} \boldsymbol{\mu}$-dominates $\mathfrak{B}$ (and write $\mathfrak{A} \succ^{\boldsymbol{\mu}}$ $\mathfrak{B})$ iff for all $\mathfrak{b} \in \mathfrak{B}$ there exists $\mathfrak{a} \in \mathfrak{A}$ such that $\mathfrak{a} \succ^{\boldsymbol{\mu}} \mathfrak{b}$. Let $S, T \in \mathbb{T}^{\boldsymbol{\mu}}$. We say $S \boldsymbol{\mu}$-dominates $T$ (and write $S \succ^{\boldsymbol{\mu}} T$ ) iff $\operatorname{supp} S \boldsymbol{\mu}$-dominates $\operatorname{supp} T$. Note that this agrees with the previous definitions for $\succ^{\boldsymbol{\mu}}$ when $S=1$ or when $S, T \in \mathfrak{J}^{\mu}$.
Remark 4.13. $S \succ T$ if and only if there exists $\boldsymbol{\mu}$ such that $S \succ^{\boldsymbol{\mu}} T$.
Remark 4.14. Use of $\succ^{\boldsymbol{\mu}}$ requires caution (at least for non-monomials), because it does not always have the properties of $\succ$. For example: $\boldsymbol{\mu}$-dominance is not preserved by multiplication. That is: $A \prec^{\mu} B$ does not imply $A S \prec^{\mu} B S$. For example, let $\boldsymbol{\mu}=\left\{x^{-1}, e^{-x}\right\}, A=x^{-2}+e^{-2 x}, B=x^{-1}+e^{-x}+x e^{-2 x}$, and $S=x^{-1}-e^{-x}$. The term $x^{-1} e^{-2 x}$ of $A S$ is not $\boldsymbol{\mu}$-dominated by any term of $B S$.

The following four propositions are proved as in multi-indices (Propositions 4.4 to 4.7 ).

Proposition 4.15. Let $\mathfrak{A}, \mathfrak{B} \subseteq \mathfrak{J}^{\boldsymbol{\mu}, \mathbf{m}}$. Then $\mathfrak{A} \succ^{\mu} \mathfrak{B}$ if and only if $\operatorname{Max}^{\mu} \mathfrak{A} \succ^{\mu}$ $\operatorname{Max}^{\mu} \mathfrak{B}$.

Proposition 4.16. Let $\mathfrak{A}, \mathfrak{B} \subseteq \mathfrak{J}^{\boldsymbol{\mu}, \mathbf{m}}$. If $\mathfrak{A} \succ^{\boldsymbol{\mu}} \mathfrak{B}$, then $\operatorname{Max}^{\boldsymbol{\mu}} \mathfrak{A}$ and $\operatorname{Max}^{\boldsymbol{\mu}} \mathfrak{B}$ are disjoint.

Proposition 4.17. Let $\mathfrak{A}_{j} \subseteq \mathfrak{J}^{\mu, \mathbf{m}}, j \in \mathbb{N}$, be an infinite sequence such that $\mathfrak{A}_{j} \succ^{\mu} \mathfrak{A}_{j+1}$ for all $j$. Then the sequence $\left(\mathfrak{A}_{j}\right)$ is point-finite

Proposition 4.18. Let $\mathfrak{A}_{i} \subseteq \mathfrak{J}^{\boldsymbol{\mu}, \mathbf{m}}$ be a point-finite family. Assume $\mathfrak{A}_{i} \succ^{\mu} \mathfrak{B}_{i}$ for all $i$. Then $\mathfrak{B}_{i} \subseteq \mathfrak{J}^{\mu, \mathbf{m}}, i \in \mathbb{N}$, and the family $\left(\mathfrak{B}_{i}\right)$ is also point-finite.
Definition 4.19. Let $\Phi$ be linear from some subspace of $\mathbb{T}^{\mu}$ to itself. Then we say $\Phi$ is $\boldsymbol{\mu}$-contractive iff $T \succ^{\boldsymbol{\mu}} \Phi(T)$ for all $T \neq 0$ in the subspace.
Definition 4.20. Let $\Phi$ be possibly non-linear from some subset $\mathbf{A}$ of $\mathbb{T}^{\mu}$ to itself. Then we say $\Phi$ is $\boldsymbol{\mu}$-contractive iff $(S-T) \succ^{\boldsymbol{\mu}}(\Phi(S)-\Phi(T))$ for all $S, T \in \mathbf{A}$ with $S \neq T$.

There is an easy way to define a linear $\boldsymbol{\mu}$-contractive map $\Phi$. If $\Phi$ is defined on all monomials $\mathfrak{g} \in \mathfrak{A} \subseteq \mathfrak{J}^{\mu, \mathbf{m}}$ and $\mathfrak{g} \succ^{\mu} \Phi(\mathfrak{g})$ for them, then the family $(\operatorname{supp} \Phi(\mathfrak{g}))$ is point-finite by Proposition 4.18, so

$$
\Phi\left(\sum c_{\mathfrak{g}} \mathfrak{g}\right)=\sum c_{\mathfrak{g}} \Phi(\mathfrak{g})
$$

$\boldsymbol{\mu}$-converges and defines $\Phi$ on the span.
Example 4.21. The set $\boldsymbol{\mu}$ of ratios is important. (In fact, this is the reason we have been paying so much attention to the ratio set $\boldsymbol{\mu}$.) We cannot simply replace " $\boldsymbol{\mu}$-small" by "small" in the definitions. Suppose $\Phi\left(x^{-j}\right)=x^{j} e^{-x}$ for all $j \in \mathbb{N}$, and $\Phi(\mathfrak{g})=\mathfrak{g} x^{-1}$ for all other monomials. Then $\mathfrak{g} \succ \Phi(\mathfrak{g})$ for all $\mathfrak{g}$. But $\Phi\left(\sum x^{-j}\right)$ evaluated termwise is not a legal transseries. Or: Define $\Phi\left(x^{-j}\right)=e^{-x}$ for all $j \in \mathbb{N}$, and $\Phi(\mathfrak{g})=\mathfrak{g} x^{-1}$ for all other monomials. Again $\mathfrak{g} \succ \Phi(\mathfrak{g})$ for all $\mathfrak{g}$, but the family $\operatorname{supp} \Phi\left(x^{-j}\right)$ is not point-finite.
Theorem 4.22 (Grid-Based Fixed-Point Theorem). (i) If $\Phi$ is linear and $\boldsymbol{\mu}$-contractive on $\mathbb{T}^{\boldsymbol{\mu}, \mathbf{m}}$, then for any $T_{0} \in \mathbb{T}^{\boldsymbol{\mu}, \mathbf{m}}$, the fixed-point equation $T=$ $\Phi(T)+T_{0}$ has a unique solution $T \in \mathbb{T}^{\boldsymbol{\mu}, \mathbf{m}}$. (ii) If $\mathbf{A} \subseteq \mathbb{T}^{\boldsymbol{\mu}, \mathbf{m}}$ is nonempty and closed (in the asymptotic topology), and nonlinear $\Phi: \mathbf{A} \rightarrow \mathbf{A}$ is $\boldsymbol{\mu}$-contractive on $\mathbf{A}$, then $T=\Phi(T)$ has a unique solution in $\mathbf{A}$. (From [11, Theorem 15]. See [37, §6.5] and [35].)

Proof. (i) follows from (ii), since if $\Phi$ is linear and $\boldsymbol{\mu}$-contractive, then $\widetilde{\Phi}$ defined by $\widetilde{\Phi}(T)=\Phi(T)+T_{0}$ is $\boldsymbol{\mu}$-contractive.
(ii) First note $\Phi$ is $\boldsymbol{\mu}$-continuous: Assume $T_{j} \xrightarrow{\boldsymbol{\mu}} T$. Then $T_{j}-T \xrightarrow{\boldsymbol{\mu}} 0$, so $\left(\operatorname{supp}\left(T_{j}-T\right)\right)$ is point-finite. But $\operatorname{supp}\left(T_{j}-T\right) \succ^{\boldsymbol{\mu}} \operatorname{supp}\left(\Phi\left(T_{j}\right)-\Phi(T)\right)$, so $\left(\operatorname{supp}\left(\Phi\left(T_{j}\right)-\Phi(T)\right)\right.$ is also point-finite by Proposition 4.18. And so $\Phi\left(T_{j}\right) \xrightarrow{\mu}$ $\Phi(T)$.

Existence: Define $T_{j+1}=\Phi\left(T_{j}\right)$. We claim $T_{j}$ is $\boldsymbol{\mu}$-convergent. The sequence $\mathfrak{A}_{j}=\operatorname{supp}\left(T_{j}-T_{j+1}\right)$ satisfies: $\mathfrak{A}_{j} \succ^{\boldsymbol{\mu}} \mathfrak{A}_{j+1}$ for all $j$, so (Proposition 4.17) $\left(\mathfrak{A}_{j}\right)$ is point-finite, which means $T_{j}-T_{j+1} \xrightarrow{\mu} 0$ and therefore
(by nonarchimedean Cauchy) $T_{j} \boldsymbol{\mu}$-converges. Difference preserves $\boldsymbol{\mu}$-limits, so the limit $T$ satisfies $\Phi(T)=T$.

Uniqueness: if $A$ and $B$ were two different solutions, then $\Phi(A)-\Phi(B)=$ $A-B$, which contradicts $\boldsymbol{\mu}$-contractivity.

Remark 4.23. The $\boldsymbol{\mu}$-dominance relation may be used to explain two of the earlier results that may have seemed poorly motivated at the time.
(a) To prove the existence of the derivative: When $\mathfrak{g}^{\prime}$ had been defined for $\mathfrak{g} \in \mathfrak{J}^{\boldsymbol{\mu , m}}$, we then showed (Proposition 3.76) that the set $\left\{\mathfrak{g}^{\prime}: \mathfrak{g} \in \mathfrak{J}^{\mu, \mathbf{m}}\right\}$ is point-finite. We could first show: given $\boldsymbol{\mu}$ there exists $\widetilde{\boldsymbol{\mu}}$ such that if $\mathfrak{m} \prec^{\mu} \mathfrak{n}$, $\mathfrak{n} \neq 1$, then $\mathfrak{m}^{\prime} \prec^{\tilde{\mu}} \mathfrak{n}^{\prime}$. Then: Given any $\mathfrak{n} \in \mathfrak{G}$, we claim that the set $\mathfrak{A}=\left\{\mathfrak{g} \in \mathfrak{J}^{\mu, \mathfrak{m}}: \mathfrak{n} \in \operatorname{supp} \mathfrak{g}^{\prime}\right\}$ is finite. If not, by Proposition 4.10 there is an infinite sequence $\mathfrak{g}_{j} \in \mathfrak{A}, \mathfrak{g}_{j} \neq 1$, with $\mathfrak{g}_{0} \succ^{\boldsymbol{\mu}} \mathfrak{g}_{1} \succ^{\boldsymbol{\mu}} \ldots$. But then $\mathfrak{g}_{0}^{\prime} \succ^{\tilde{\mu}} \mathfrak{g}_{1}^{\prime} \succ^{\tilde{\mu}} \cdots$, so $\left\{\operatorname{supp} \mathfrak{g}_{j}^{\prime}: j \in \mathbb{N}\right\}$ is point-finite by Proposition 4.17, contradicting the assumption that $\mathfrak{A}$ is infinite.
(b) To prove the existence of the composition $T \circ S$ : When $\mathfrak{g} \circ S$ had been defined for $\mathfrak{g} \in \mathfrak{J}^{\mu, \mathrm{m}}$, we then showed (Proposition 3.95) that the set $\left\{\mathfrak{g} \circ S: \mathfrak{g} \in \mathfrak{J}^{\boldsymbol{\mu}, \mathbf{m}}\right\}$ is point-finite. We could first show: given $\boldsymbol{\mu}$ and $S$, there exists $\tilde{\boldsymbol{\mu}}$ such that if $\mathfrak{m} \succ^{\boldsymbol{\mu}} \mathfrak{n}$, then $\mathfrak{m} \circ S \succ^{\tilde{\boldsymbol{\mu}}} \mathfrak{n} \circ S$. Then: Given any $\mathfrak{n} \in \mathfrak{G}$, we claim that the set $\mathfrak{A}=\left\{\mathfrak{g} \in \mathfrak{J}^{\boldsymbol{\mu , m}}: \mathfrak{n} \in \operatorname{supp}(\mathfrak{g} \circ S)\right\}$ is finite. If not, by Proposition 4.10 there is an infinite sequence $\mathfrak{g}_{j} \in \mathfrak{A}$ with $\mathfrak{g}_{0} \succ^{\boldsymbol{\mu}} \mathfrak{g}_{1} \succ^{\mu} \ldots$. But then $\mathfrak{g}_{0} \circ S \succ^{\tilde{\mu}} \mathfrak{g}_{1} \circ S \succ^{\tilde{\mu}} \cdots$, so $\left\{\operatorname{supp}\left(\mathfrak{g}_{j} \circ S\right): j \in \mathbb{N}\right\}$ is point-finite by Proposition 4.17, contradicting the assumption that $\mathfrak{A}$ is infinite.

## Integration.

In elementary calculus courses, we find that certain integrals can be evaluated using reduction formulas. For example $\int x^{n} e^{x} d x$, when integrated by parts, yelds an integral of the same form, but with exponent $n$ reduced by 1. So if we repeat this until the exponent is zero, we have our integral. But of course this does not work when the exponent is not in $\mathbb{N}$. We can try it, and get an infinite series:

$$
\begin{aligned}
\int x^{a} e^{x} d x & =x^{a} e^{x}-a \int x^{a-1} e^{x} d x \\
& =x^{a} e^{x}-a x^{n-1} e^{x}+a(a-1) \int x^{a-2} e^{x} d x=\cdots \\
& =\sum_{j=0}^{\infty}(-1)^{j} a(a-1)(a-2) \cdots(a-j+1) x^{a-j} e^{x}
\end{aligned}
$$

$$
=\sum_{j=0}^{\infty} \frac{\Gamma(j-a)}{\Gamma(-a)} x^{a-j} e^{x}
$$

but this series converges for no $x$. However, it is still a transseries solution to the problem.

Proposition 4.24. Let $a, b, c \in \mathbb{R}, c>0, b \neq 0$. Then the transseries

$$
T=\sum_{j=1}^{\infty} \frac{\Gamma\left(j-\frac{a+1}{c}\right)}{\Gamma\left(1-\frac{a+1}{c}\right) c b^{j}} x^{a+1-j c} e^{b x^{c}}
$$

has derivative $T^{\prime}=x^{a} e^{b x^{c}} .(I f(a+1) / c$ is a positive integer, then $T$ should be a finite sum.)

Problem 4.25. More generally: if $b \in \mathbb{R}$ and $L \in \mathbb{T}_{0}$ is purely large, can you use the same method to show that there is $T \in \mathbb{T}_{1}$ with $T^{\prime}=x^{b} e^{L}$ ?

## The General Integral

Every transseries in $\mathbb{R} \mathbb{L} x \mathbb{W}$ has an integral (an antiderivative). We will give a complete proof following the hint in [13, 4.10e.1]. (See also [18, Prop. 7.22].) This is an example where we convert the problem to a log-free case to apply the contraction argument. The general integration problem (Theorem 4.29) is reduced to one (Proposition 4.26) where contraction can be easily applied.

Proposition 4.26. Let $T \in \mathbb{T}_{\bullet}$ with $T \succ 1$. Then there is $S \in \mathbb{T}_{\bullet}$ with $S^{\prime}=e^{T}$.

Proof. Either $T$ is positive or negative. We will do the positive case, the negative one is similar (and it turns out the iterative formulas are the same). If

$$
S=\frac{e^{T}}{T^{\prime}}(1+U)
$$

where $U$ satisfies

$$
U=\frac{T^{\prime \prime}}{\left(T^{\prime}\right)^{2}}+\frac{T^{\prime \prime}}{\left(T^{\prime}\right)^{2}} U-\frac{U^{\prime}}{T^{\prime}}
$$

then it is a computation to see that $S^{\prime}=e^{T}$. So it suffices to exhibit an appropriate $\boldsymbol{\mu}$ and show that the linear map $\Phi: \mathbb{T}^{\boldsymbol{\mu}, \mathbf{0}} \rightarrow \mathbb{T}^{\boldsymbol{\mu}, \mathbf{0}}$ defined by

$$
\Phi(U)=\frac{T^{\prime \prime}}{\left(T^{\prime}\right)^{2}} U-\frac{U^{\prime}}{T^{\prime}}
$$

is $\boldsymbol{\mu}$-contractive, then apply Theorem $4.22(\mathrm{i})$.

Say $T$ is of exact height $N$, so $e^{T}$ is of exact height $N+1$. By Proposition 3.82, $T^{\prime \prime} \prec\left(T^{\prime}\right)^{2}$ and $x T^{\prime} \succ 1$. So $T^{\prime \prime} /\left(T^{\prime}\right)^{2}$ and $1 /\left(x T^{\prime}\right)$ are small. Let $\tilde{\boldsymbol{\mu}}$ be the least set of ratios including $x^{-1}$, ratios generating $T$, the inversion addendum for $T^{\prime}$, the smallness addenda for $T^{\prime \prime} /\left(T^{\prime}\right)^{2}$ and $1 /\left(x T^{\prime}\right)$, and is hereditary. Then, for each $\mu_{i}=x^{-b_{i}} e^{-L_{i}}$ in $\widetilde{\boldsymbol{\mu}}$ (finitely many of them), since $L_{i}$ has lower height than $T^{\prime}$, we have $L_{i}^{\prime} / T^{\prime} \prec 1$. Add smallness addenda for all of these, call the result $\boldsymbol{\mu}$. Note $\operatorname{lsupp} \boldsymbol{\mu}=\operatorname{lsupp} \tilde{\boldsymbol{\mu}}$, so we don't have to repeat this last step.

By Proposition 3.86 all ratios in $\boldsymbol{\mu}$ are (at most) of height $N$. And all derivatives $T^{\prime}, T^{\prime \prime}$ belong to $\mathbb{T}^{\mu}$. The function $\Phi$ maps $\mathbb{T}^{\mu}$ into itself.

Since $\Phi$ is linear, we just have to check for monomials $\mathfrak{g} \in \mathfrak{J}^{\mu, \mathbf{0}}$ that $\mathfrak{g} \succ^{\mu} \Phi(\mathfrak{g})$. Now $T^{\prime \prime} /\left(T^{\prime}\right)^{2}$ is $\boldsymbol{\mu}$-small so $\mathfrak{g} \succ^{\boldsymbol{\mu}}\left(T^{\prime \prime} /\left(T^{\prime}\right)^{2}\right) \mathfrak{g}$. For the second term: If $\mathfrak{g}=\boldsymbol{\mu}^{\mathbf{k}}=x^{b} e^{L}$, then

$$
\frac{\mathfrak{g}^{\prime}}{T^{\prime}}=\frac{b x^{b-1} e^{L}+L^{\prime} x^{b} e^{L}}{T^{\prime}}=\frac{b x^{-1}+L^{\prime}}{T^{\prime}} \mathfrak{g}=\frac{b}{x T^{\prime}} \mathfrak{g}+\frac{L^{\prime}}{T^{\prime}} \mathfrak{g}
$$

But $1 / x T^{\prime} \prec^{\boldsymbol{\mu}} 1$ so $\mathfrak{g} \succ^{\boldsymbol{\mu}}\left(b /\left(x T^{\prime}\right)\right) \mathfrak{g}$. And $L^{\prime} / T^{\prime} \prec^{\boldsymbol{\mu}} 1$ so $\mathfrak{g} \succ^{\boldsymbol{\mu}}\left(L^{\prime} / T^{\prime}\right) \mathfrak{g}$.
Definition 4.27. We say $x^{b} e^{L} \in \mathfrak{G}_{\bullet}$ is power-free iff $b=0$. We say $T \in \mathbb{T}_{\bullet}$ is power-free iff all transmonomials in supp $T$ are power-free.

Since $\left(x^{b} e^{L}\right) \circ \exp =e^{b x} e^{L \circ e x p}$, it follows that all $T \in \mathbb{T}_{\bullet,-1}$ are power-free.
Proposition 4.28. Let $T \in \mathbb{T}_{\bullet}$ be a power-free transseries. Then there is $S \in \mathbb{T}_{\bullet}$ with $S^{\prime}=T$.

Proof. For monomials $\mathfrak{g}=e^{L} \in \operatorname{supp} T$ with large $L \in \mathbb{T}^{\mu}$, write $\mathcal{P}(\mathfrak{g})$ for the transseries constructed in Proposition 4.26 with $\mathcal{P}(\mathfrak{g})^{\prime}=\mathfrak{g}$. Then we must show that the family $(\operatorname{supp} \mathcal{P}(\mathfrak{g}))$ is point-finite, so we can define $\mathcal{P}\left(\sum c_{\mathfrak{g}} \mathfrak{g}\right)=\sum c_{\mathfrak{g}} \mathcal{P}(\mathfrak{g})$. For large $L$ we have $x L^{\prime} \succ 1$ (Proposition 3.82). Thus, the formula

$$
\frac{\mathcal{P}\left(e^{L}\right)}{x}=\frac{e^{L}}{x L^{\prime}}(1+U)
$$

shows that $e^{L} \tilde{\boldsymbol{\mu}}$-dominates $\mathcal{P}\left(e^{L}\right) / x$ for an appropriate $\tilde{\boldsymbol{\mu}}$. So the family $\operatorname{supp}\left(\mathcal{P}\left(e^{L}\right) / x\right)$ is point-finite and thus the family $\operatorname{supp} \mathcal{P}\left(e^{L}\right)$ is point-finite.

Theorem 4.29. Let $A \in \mathbb{R} \llbracket x \rrbracket$. Then there exists $B \in \mathbb{R} \llbracket x \rrbracket$ with $B^{\prime}=A$.
Proof. Say $A \in \mathbb{T}_{\bullet M}$. Then $A=T_{1} \circ \log _{M+1}$, where $T_{1} \in \mathbb{T}_{\bullet,-1}$. Let

$$
T=T_{1} \cdot \exp _{M+1} \cdot \exp _{M} \cdots \exp _{2} \cdot \exp _{1}
$$

Now $T$ is power-free, so by Proposition 4.28, there exists $S \in \mathbb{T} \bullet$ with $S^{\prime}=T$. Then let $B=S \circ \log _{M+1}$ and check that $B^{\prime}=A$. Note that $B \in \mathbb{T}_{\bullet, M+1}$.

## An Integral

Problem 4.30. Compute the integral

$$
\int e^{e^{e^{x}}} d x
$$

using the method of Proposition 4.26.
We first display the ratio set to be used, and derivatives:

$$
\begin{array}{lll}
\mu_{1}=x^{-1}, & \mu_{1}^{\prime}=-\mu_{1}^{2}, & L_{1}=\mu_{1}^{-1}=x \\
\mu_{2}=e^{-x}=e^{-L_{1}}, & \mu_{2}^{\prime}=-\mu_{2}, & L_{2}=\mu_{2}^{-1}=e^{x} \\
\mu_{3}=e^{-e^{x}}=e^{-L_{2}}, & \mu_{3}^{\prime}=-\mu_{2}^{-1} \mu_{3}, & L_{3}=\mu_{3}^{-1}=e^{e^{x}} \\
\mu_{4}=e^{-e^{e^{x}}}=e^{-L_{3}}, & \mu_{4}^{\prime}=-\mu_{2}^{-1} \mu_{3}^{-1} \mu_{4}, & L_{4}=\mu_{4}^{-1}=e^{e^{e^{x}}}
\end{array}
$$

The integral should have the form $\left(L_{4} / L_{3}^{\prime}\right)(1+U)$, where $U$ satisfies

$$
U=\frac{L_{3}^{\prime \prime}}{\left(L_{3}^{\prime}\right)^{2}}+\frac{L_{3}^{\prime \prime}}{\left(L_{3}^{\prime}\right)^{2}} U-\frac{U^{\prime}}{L_{3}^{\prime}}=\left(\mu_{3}+\mu_{2} \mu_{3}\right)+\left(\mu_{3}+\mu_{2} \mu_{3}\right) U-\mu_{2} \mu_{3} U^{\prime}
$$

To solve this, we should iterate $U_{n+1}=\left(\mu_{3}+\mu_{2} \mu_{3}\right)+\Phi\left(U_{n}\right)$ where $\Phi(Y)=$ $\left(\mu_{3}+\mu_{2} \mu_{3}\right) Y-\mu_{2} \mu_{3} Y^{\prime}$. Starting with $U_{0}=0$, we get

$$
\begin{aligned}
& U_{1}=\left(1+\mu_{2}\right) \mu_{3} \\
& U_{2}=\left(1+\mu_{2}\right) \mu_{3}+\left(2+3 \mu_{2}+2 \mu_{2}^{2}\right) \mu_{3}^{2} \\
& U_{3}=\left(1+\mu_{2}\right) \mu_{3}+\left(2+3 \mu_{2}+2 \mu_{2}^{2}\right) \mu_{3}^{2}+\left(6+11 \mu_{2}+12 \mu_{2}^{2}+6 \mu_{2}^{3}\right) \mu_{3}^{3}
\end{aligned}
$$

each step producing one higher power of $\mu_{3}$ and preserving all of the existing terms. Once we have the limit $U$, we add 1 and multiply by $L_{4} / L_{3}^{\prime}=$ $e^{e^{e^{x}}} /\left(e^{x} e^{e^{x}}\right)$. The result is

$$
\int e^{e^{e^{x}}} d x=e^{e^{e^{x}}} \sum_{j=1}^{\infty} e^{-j e^{x}}\left(\sum_{k=1}^{j} e^{-k x} c_{j, k}\right)
$$

The coefficients $c_{j, k}$ (namely, $1 ; 1,1 ; 2,3,2 ; 6,11,12,6 ; \cdots$ ) are related to Stirling numbers of the second kind.

Similarly, we may compute

$$
\int e^{k_{2} x} e^{k_{3} e^{x}} e^{k_{4} e^{e^{x}}} d x=e^{k_{2} x} e^{k_{3} e^{x}} e^{k_{4} e^{e^{x}}} \sum_{j=1}^{\infty} e^{-j e^{x}}\left(\sum_{k=1}^{j} e^{-k x} c_{j, k}\right)
$$

for some coefficients $c_{j, k}$ depending on $k_{2}, k_{3}, k_{4}$.

## A Differential Equation

Problem 4.31. Solve the Riccati equation

$$
\begin{equation*}
Y^{\prime}=\frac{x-x^{2}}{x^{2}-x+1} Y+Y^{2} \tag{*}
\end{equation*}
$$

This is a differential equation where the solution can be written in closed form. (At least if you consider an integral to be closed form.) But it will illustrate some things to watch out for when computing transseries solutions. The same things can happen in cases where solutions are not known in closed form.

If we are not careful, we may come up with a series

$$
S(x)=e^{-x}\left(1-\frac{1}{x}+\frac{1}{3 x^{3}}+\frac{1}{6 x^{4}}-\frac{1}{10 x^{5}}-\frac{8}{45 x^{6}}-\frac{1}{18 x^{7}}+\frac{11}{120 x^{8}}+\cdots\right)
$$

and claim it is a solution. If we plug this series in for $Y$, then the two sides of the differential equation agree to all orders. That is, if we compute $S$ up to $\mathrm{O}\left(e^{-x} x^{-1000}\right)$, and plug it in, then the two sides agree up to $\mathrm{O}\left(e^{-x} x^{-1000}\right)$. But in fact, $S$ is not the transseries solution of (*). The two sides are not equal-their difference is just far smaller than all terms of the series $S$. The difference has order $e^{-2 x}$. In hindsight, this should be clear, because of the $Y^{2}$ term in $(*)$. If $Y$ has a term $e^{-x}$ in its expansion, then $Y^{2}$ will have a term $e^{-2 x}$. When $S$ is substituted into $(*)$, the term $e^{-2 x}$ appears on the right side but not the left.

In fact, $S(x)$ is a solution of $(*)$ without the $Y^{2}$ term.
According to Maple, the actual solution is $Y=c S(x) /\left(1-c \int S(x) d x\right)$, where $c$ is an arbitrary constant and

$$
S(x)=\exp \left[-x+\frac{2}{\sqrt{3}} \arctan \left(\frac{2 x-1}{\sqrt{3}}\right)\right]
$$

The exponent in this $S(x)$ is $-x$ plus constant plus small, so $S$ can be written as a series. It is (except for the constant factor) the series $S(x)$ given above.

Now the integral of $S$ can be done (using Proposition 4.24), then division carried out as usual. The general solution of $(*)$ is:

$$
\begin{aligned}
& c e^{-x}\left(1-\frac{1}{x}+\frac{1}{3 x^{3}}+\frac{1}{6 x^{4}}-\frac{1}{10 x^{5}}-\frac{8}{45 x^{6}}+\cdots\right) \\
+ & c^{2} e^{-2 x}\left(-1+\frac{2}{x}-\frac{-2}{x^{2}}+\frac{7}{3 x^{3}}-\frac{20}{3 x^{4}}+\frac{388}{15 x^{5}}-\frac{5578}{45 x^{6}}+\cdots\right)
\end{aligned}
$$

$$
\begin{aligned}
& +c^{3} e^{-3 x}\left(1-\frac{3}{x}+\frac{5}{x^{2}}-\frac{8}{x^{3}}+\frac{39}{2 x^{4}}-\frac{693}{10 x^{5}}+\frac{3159}{10 x^{6}}+\cdots\right) \\
& +c^{4} e^{-4 x}\left(-1+\frac{4}{x}-\frac{9}{x^{2}}+\frac{53}{3 x^{3}}-\frac{128}{3 x^{4}}+\frac{707}{5 x^{5}}-\frac{27442}{45 x^{6}}+\cdots\right) \\
& +c^{5} e^{-5 x}\left(1-\frac{5}{x}+\frac{14}{x^{2}}-\frac{97}{3 x^{3}}+\frac{487}{6 x^{4}}-\frac{1549}{6 x^{5}}+\frac{9509}{9 x^{6}}+\cdots\right) \\
& +c^{6} e^{-6 x}\left(-1+\frac{6}{x}-\frac{20}{x^{2}}+\frac{53}{x^{3}}-\frac{141}{x^{4}}+\frac{2208}{5 x^{5}}-\frac{8648}{5 x^{6}}+\cdots\right) \\
& +\cdots
\end{aligned}
$$

Transseries solutions to simple problems can have support of transfinite order type!

The transseries solution to differential equation $(*)$ can be found without using a known closed form. The generic method would reduce to height zero (by taking logarithms of the unknown $Y$ ) then solve as a contractive map.

There is another comment on doing these computations with a computer algebra system. Carrying out the division indicated above, for example, is not trivial. If I write the two series to many terms, divide, then tell Maple to write it as a series (using the MultiSeries package, series ( $A / B, x=i n f i n i t y, 15$ )), I get only the first row of the result above. Admitedly, there is a big-O term at the end, and all terms in the subsequent rows are far smaller than that, but it is not what we want here.

We want to discard not terms that are merely small, but terms that are $\boldsymbol{\mu}$-small for a relevant $\boldsymbol{\mu}$. So this computation can better be done using a grid. Choose a finite ratio set-in this case I used $\mu_{1}=x^{-1}, \mu_{2}=e^{-x}$. We write the two series in terms of these ratios, then expand the quotient as a series in the two variables $\mu_{1}, \mu_{2}$. Now we can control which terms are kept. Delete terms not merely when they are small, but when they are $\boldsymbol{\mu}$-small. The series above has all terms $\boldsymbol{\mu}^{\mathbf{k}}$ with $\mathbf{k} \leq(6,6)$.

I used this same grid method for the computations in Problem 4.30. That is the reason I started there by displaying the required ratio set and derivatives.

## Factor a Differential Operator

Problem 4.32. Factor the differential operator

$$
\partial^{2}+x \partial+I=(\partial-\alpha(x) I)(\partial-\beta(x) I)
$$

where $\alpha(x)$ and $\beta(x)$ are transseries.

Why don't you do it? My answer looks like this:

$$
\begin{aligned}
& \alpha(x)=-x+\frac{1}{x}+\frac{2}{x^{3}}+\frac{10}{x^{5}}+\frac{74}{x^{7}}+\frac{706}{x^{9}}+\frac{8162}{x^{11}}+\frac{110410}{x^{13}}+\mathrm{O}\left(x^{-15}\right) \\
& \beta(x)=-\frac{1}{x}-\frac{2}{x^{3}}-\frac{10}{x^{5}}-\frac{74}{x^{7}}-\frac{706}{x^{9}}-\frac{8162}{x^{11}}-\frac{110410}{x^{13}}+\mathrm{O}\left(x^{-15}\right)
\end{aligned}
$$

Are the coefficients Sloane A000698 [56]? I am told that these series are divergent, and that can be proved by considering "Stokes directions" in the complex plane - another interesting topic beyond the scope of this paper. Elementary functions have convergent transseries [17, Cor. 5.5], so $\alpha(x)$ and $\beta(x)$ (even the genuine functions obtained by Écalle-Borel summation) are not elementary functions.

## Increasing and Decreasing

I have (as part of the sales pitch) tried to show that the reasoning required for the theory of transseries is easy, although perhaps sometimes tedious. But, in fact, I think there are situations - dealing with composition-that are not as easy.
Problem 4.33. Let $T, A, B \in \mathbb{T}$. Assume $A, B$ are large and positive. Prove or disprove: if $T^{\prime}>0$ and $A<B$, then $T \circ A<T \circ B$.

## 5 Additional Remarks.

If (as I claim) the system $\mathbb{R} \llbracket x \rrbracket$ of transseries is an elementary and fundamental object, then perhaps it is only natural that there are variants in the formulation and definitions used. For example [37] the construction can proceed by first adding logarithms, and then adding exponentials. For an exercise, see if you can carry that out yourself in such a way that the end result is the same system of transseries as constructed above. I prefer the approach shown here, since I view the "log-free" calculations as fundamental.

There is a possibility $[1,23,34,43,54]$ to allow well ordered supports instead of just the grids $\mathfrak{J}^{\mu, \mathbf{m}}$. These are called well-based transseries. (Perhaps we use the alternate notation $\mathbb{R}[[\mathfrak{M}]]$ for the well-based Hahn field and the new notation $\mathbb{R} \llbracket \mathfrak{M} \rrbracket$ for the grid-based subfield.) The set of well-based transseries forms a strictly larger system than the grid-based transseries, but with most of the same properties. Which of these two is to be preferred may be still open to debate. In this paper we have used the grid-based approach because:
(i) The finite ratio set is conducive to computer calculations.
(ii) Problems from analysis almost always have solutions in this smaller system.
(iii) Some proofs and formulations of definitions are simpler in one system than in the other.
(iv) Perhaps (?) the analysis used for Écalle-Borel convergence can be applied only to grid-based series.
(v) [42] In the well-based case, the domain of exp cannot be all of $\mathbb{R}[[\mathfrak{M}]]$.
(vi) [31] The grid-based ordered set $\mathbb{R} \llbracket \mathfrak{M} \rrbracket$ is a "Borel order," but the wellbased ordered set $\mathbb{R}[[\mathfrak{M}]]$ is not.

What constitutes desirable properties of "fields of transseries" has been explored axiomatically by Schmeling [54]. Kuhlmann and Tressl [44] compare different constructions for ordered fields of generalized power series with logarithm and exponential; the grid-based transseries we have used here might be thought of as the smallest such system (without restricting the coefficients).

Just as the real number system $\mathbb{R}$ is extended to the complex numbers $\mathbb{C}$, there are ways to extend the system of real transseries to allow for complex numbers. The simplest uses the same group $\mathfrak{G}_{\bullet .}$ of monomials, but then takes complex coefficients to form $\mathbb{C}[\mathfrak{G} \bullet \bullet \rrbracket]=\mathbb{C} \llbracket x \rrbracket]$. For example, the fifthdegree equation in Problem 4.1 has five solutions in $\mathbb{C}[\mathfrak{G} \bullet \bullet]$. But this still won't give us oscillatory functions, such as solutions to the differential equation $Y^{\prime \prime}+Y=0$. There is a way [37, Section 7.7] to define oscillating transseries. These are finite sums

$$
\sum_{j=1}^{n} \alpha_{j} e^{i \psi_{j}}
$$

with amplitudes $\alpha_{j} \in \mathbb{C} \llbracket x \rrbracket$ and purely large phases $\psi_{j} \in \mathbb{R} \llbracket x \rrbracket$. And van der Hoeven [36] considers defining complex transseries using the same method as we used for real transseries, where the required orderings are done in terms of sectors in the complex plane.

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