# Boolean Witt vectors and an integral Edrei-Thoma theorem 

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#### Abstract

A subtraction-free definition of the big Witt vector construction was recently given by the first author. This allows one to define the big Witt vectors of any semiring. Here we give an explicit combinatorial description of the big Witt vectors of the Boolean semiring. We do the same for two variants of the big Witt vector construction: the Schur Witt vectors and the $p$-typical Witt vectors. We use this to give an explicit description of the Schur Witt vectors of the natural numbers, which can be viewed as the classification of totally positive power series with integral coefficients, first obtained by Davydov. We also determine the cardinalities of some Witt vector algebras with entries in various concrete arithmetic semirings.


Keywords Witt vector • Semiring • Symmetric function • Schur positivity • Total positivity • Boolean algebra

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## 1 Introduction

Recently in [5], the first author gave an extension of the big Witt vector functor to a functor $W$ on (commutative) semirings-in other words, a definition of Witt vectors that does not require subtraction. Witt vectors are usually thought of as being related to the $p$-adic numbers or, more broadly, the finite adeles. For example, the (relative) de Rham-Witt complex, in its single-prime form, computes crystalline cohomology; in its multiple-prime form, it formally unifies the crystalline cohomologies at all primes. (See, e.g., $[3,6,12-14,17]$.) But the extension of the Witt construction to semirings shows that positivity also plays a role, and thus so does the archimedean prime. One striking expression of this is that $W(\mathbb{N})$ is countable and has no zero divisors, and therefore can connect to analytic mathematics in ways that $W(\mathbb{Z})$, which is pro-discrete and has many zero divisors, cannot. For instance, there is a canonical injective homomorphism from $W(\mathbb{N})$ to the ring of entire functions on $\mathbb{C}$. In [4], it was argued that it is reasonable to think of $W(\mathbb{Z})$ as $\mathbb{Z} \otimes_{\mathbb{F}_{1}} \mathbb{Z}$, where $\mathbb{F}_{1}$ is the hypothetical field with one element; so it is notable that $W(\mathbb{N})$, a more fundamental alternative, has some analytic nature.

For reasons like these, one would like to understand better the range of positivity phenomena that can occur in Witt vectors. Here a special place is held by the Boolean semiring $\mathbb{B}=\{0,1\}$, with $1+1=1$. This is because, aside from the fields, it is the unique nonzero semiring which has no quotients other than itself and 0. (See Golan [11, p. 87].) One might imagine $\mathbb{B}$ as the residue field of $\mathbb{N}$ at the infinite prime. In fact, as we show in (7.9), any semiring $A$ not containing -1 necessarily admits a map to $\mathbb{B}$, much as any ring not containing $1 / p$ admits a map to some field of characteristic $p$. In particular, there is a map $W(A) \rightarrow W(\mathbb{B})$, and hence structure related to Boolean Witt vectors is necessarily present in the Witt vectors of any semiring that is not a ring. So not only is an understanding of $W(\mathbb{B})$ helpful in understanding positivity in Witt vectors, it is all but required.

The purpose of this paper is then to give an explicit description of $W(\mathbb{B})$. From the point of view above, this would be analogous to the description of Witt vectors of fields of characteristic $p$ in concrete $p$-adic terms (which incidentally explains why $p$-adic phenomena appear in the Witt vectors of any ring not containing $1 / p$ ). We also give an explicit description $W^{\text {Sch }}(\mathbb{B})$, where $W^{\text {Sch }}$ is the Schur Witt vector functor of [5].

In order to state our main results precisely, we need to recall the basic definitions of $W$ and $W^{\text {Sch }}$ for semirings. For details, see [5]. Let $\Lambda_{\mathbb{Z}}$ denote the usual ring of symmetric functions in the formal variables $x_{1}, x_{2}, \ldots$ It is the inverse limit in the category of graded rings of the rings $\mathbb{Z}\left[x_{1}, \ldots, x_{n}\right]^{S_{n}}$ of symmetric polynomials, where the transition maps are given by $f\left(x_{1}, \ldots, x_{n}\right) \mapsto f\left(x_{1}, \ldots, x_{n-1}, 0\right)$, and where each variable $x_{i}$ has degree 1 . Then $\Lambda_{\mathbb{Z}}$ has two standard $\mathbb{Z}$-linear bases: the monomial symmetric functions $\left(m_{\lambda}\right)_{\lambda \in P}$ and the Schur symmetric functions $\left(s_{\lambda}\right)_{\lambda \in P}$, both of which are indexed by the set $P$ of partitions. (We refer the reader to chapter 1 of Macdonald's book [18] for the basics on symmetric functions.) If we write $\Lambda_{\mathbb{N}}=$ $\bigoplus_{\lambda \in P} \mathbb{N} m_{\lambda}$ and $\Lambda^{\text {Sch }}=\bigoplus_{\lambda \in P} \mathbb{N} s_{\lambda}$, where $\mathbb{N}=\{0,1, \ldots\}$, then by some standard positivity facts in the theory of symmetric functions, both are subsemirings of $\Lambda_{\mathbb{Z}}$ and we have $\Lambda^{\text {Sch }} \subseteq \Lambda_{\mathbb{N}}$. For any semiring $A$, we then define

$$
W(A)=\operatorname{Hom}_{\text {semiring }}\left(\Lambda_{\mathbb{N}}, A\right), \quad W^{\text {Sch }}(A)=\operatorname{Hom}_{\text {semiring }}\left(\Lambda^{\text {Sch }}, A\right)
$$

We call these the sets of big Witt vectors and Schur Witt vectors with entries in $A$. When $A$ is a ring, both agree with the usual ring of big Witt vectors defined by $\operatorname{Hom}_{\operatorname{Alg}_{\mathbb{Z}}}\left(\Lambda_{\mathbb{Z}}, A\right)$ in the sense that any map from $\Lambda^{\text {Sch }}$ or $\Lambda_{\mathbb{N}}$ to $A$ extends uniquely to $\Lambda_{\mathbb{Z}}$. The semiring structures on $W(A)$ and $W^{\text {Sch }}(A)$ are the unique functorial ones that agree with the usual Witt vector ring structure when $A$ is a ring. Their existence follows from some further standard positivity results in the theory of symmetric functions.

Our main result is the following combinatorial description of $W(\mathbb{B})$ and $W^{\text {Sch }}(\mathbb{B})$ :
Theorem 1.1 (1) There is a bijection

$$
\mathbb{N}^{2} \cup\{\infty\} \xrightarrow{\sim} W(\mathbb{B})
$$

sending $(x, y)$ to the unique Witt vector satisfying, for all $\lambda \in P$,

$$
m_{\lambda} \longmapsto \begin{cases}0 & \text { if } \lambda_{x+1} \geq y+1  \tag{1.1.1}\\ 1 & \text { otherwise },\end{cases}
$$

and sending $\infty$ to the unique Witt vector satisfying $m_{\lambda} \mapsto 1$ for all $\lambda \in P$.

## (2)

In these terms, the semiring structure on $W(\mathbb{B})$ is given by the laws

$$
\begin{aligned}
(x, y)+\left(x^{\prime}, y^{\prime}\right) & =\left(x+x^{\prime}, \max \left\{y, y^{\prime}\right\}\right), \quad(0, y) \cdot\left(0, y^{\prime}\right)=\left(0, \min \left\{y, y^{\prime}\right\}\right) \\
0=(0,0), \quad 1 & =(1,0), \quad \infty+z=\infty=\infty \cdot \infty, \quad \infty \cdot(0, y)=(0, y)
\end{aligned}
$$

and the pointwise partial order on $W(\mathbb{B})$ (with $0<1$ ) corresponds to the usual, componentwise partial order on $\mathbb{N}^{2}$ with terminal element $\infty$.
(3) The obvious analogues of parts (1) and (2) for $W^{\text {Sch }}(\mathbb{B})$ and the Schur basis $\left(s_{\lambda}\right)_{\lambda \in P}$ also hold except that the semiring structure on $W^{\text {Sch }}(\mathbb{B})$ is given by the laws

$$
\begin{aligned}
(x, y)+\left(x^{\prime}, y^{\prime}\right) & =\left(x+x^{\prime}, y+y^{\prime}\right), \quad(0, y) \cdot\left(0, y^{\prime}\right)=\left(y y^{\prime}, 0\right) \\
0=(0,0), \quad 1 & =(1,0), \quad \infty+z=\infty=\infty \cdot \infty, \quad \infty \cdot(0, y)=\infty
\end{aligned}
$$

The proof is in Sects. 4 and 5. We also describe the $\Lambda_{\mathbb{N}}$-semiring structure on $W(\mathbb{B})$ and the $\Lambda^{\text {Sch }}$-semiring structure on $W^{\text {Sch }}(\mathbb{B})$. Observe that both kinds of Boolean Witt vectors form countably infinite sets. For rings, this never happens.

At this point, it should be mentioned that Connes has also proposed a notion of Witt vectors with entries in certain $\mathbb{B}$-algebras [7]. See also Connes-Consani [8]. His construction is an analogue for $p=1$ of the $p$-typical Witt vector ring $W_{(p)}(K)$ with entries in a perfect field $K$ of characteristic $p$. Both the characteristic of the input ring $K$ and the typicality of the Witt construction are imagined to be 1 . In these terms, our Boolean Witt vectors have only the characteristic of the input set to 1 . The Witt vector functors themselves are the usual big (and, below, p-typical) functors, although
extended to semirings. It would be interesting to know if there are any relations between the two approaches.

In Sects. 6 and 7, we use the combinatorial description of Boolean Witt vectors above to determine $W^{\text {Sch }}(\mathbb{N})$. Our approach is by analyzing the diagram

where $\mathbb{R}_{+}=\{x \in \mathbb{R} \mid x \geq 0\}$. The maps $\sigma$ send any Witt vector $a$ to the series $\sum_{n \geq 0} a\left(e_{n}\right) t^{n}$, where

$$
e_{n}=\sum_{i_{1}<\cdots<i_{n}} x_{i_{1}} \cdots x_{i_{n}} \in \Lambda^{\mathrm{Sch}}
$$

is the $n$th elementary symmetric function, and the horizontal maps are the functorial ones. (The semiring map $\mathbb{R}_{+} \rightarrow \mathbb{B}$ is the unique one. It sends 0 to 0 and all $x>0$ to 1.) The proof is a short argument linking the combinatorial description of $W^{\mathrm{Sch}}(\mathbb{B})$ given above with an analytical description of the real column in (1.1.2) given by the Edrei-Thoma theorem from the theory of total positivity, as recalled in Sect. 6. The result is the following:

Theorem 1.2 The map $\sigma$ is a bijection from $W^{S c h}(\mathbb{N})$ to the set of integral series $f(t) \in 1+t \mathbb{Z}[[t]]$ of the form

$$
f(t)=\frac{g(t)}{h(t)},
$$

for some (necessarily unique) integral polynomials $g(t), h(t) \in 1+t \mathbb{Z}[t]$ such that all the complex roots of $g(t)$ are negative real numbers and all those of $h(t)$ are positive real numbers.

This theorem is obviously equivalent to the following one, which was apparently first observed by Davydov [9] (end of the proof of theorem 1) and deserves to be more widely known:

Theorem 1.3 An integral series $f(t) \in 1+t \mathbb{Z}[[t]]$ is totally positive if and only if it is of the form

$$
f(t)=\frac{g(t)}{h(t)},
$$

for some integral polynomials $g(t), h(t) \in 1+t \mathbb{Z}[t]$ such that all the complex roots of $g(t)$ are negative real numbers and all those of $h(t)$ are positive real numbers.

Davydov's argument also uses the Edrei-Thoma theorem, but where we use the combinatorial calculation of $W^{\text {Sch }}(\mathbb{B})$, he uses a theorem of Salem's [21].

Let us now return to Theorem (1.1). Observe that $W(\mathbb{B})$ and $W^{\text {Sch }}(\mathbb{B})$ are isomorphic as partially ordered sets, since both agree with $\mathbb{N}^{2} \cup\{\infty\}$, but that they are not isomorphic as semirings. The reason they both agree with $\mathbb{N}^{2} \cup\{\infty\}$ as sets is ultimately that nonzero rectangular partitions, which are indexed by $\mathbb{N}^{2}$, have a certain primality role in the combinatorics of both monomial and Schur symmetric functions. In fact, most of the arguments for $W(\mathbb{B})$ and for $W^{\text {Sch }}(\mathbb{B})$ are remarkably parallel. It would be interesting to have a conceptual reason why this should be so.

This surprising coincidence continues with the $p$-typical Witt vectors, which are the subject of Sect. 8 . Let $p$ be a prime number, and let $W_{(p)}$ denote the $p$-typical Witt functor for semirings, as defined in [5] and recalled in (8.1) below.

Theorem 1.4 $W_{(p)}(\mathbb{B})$ is canonically isomorphic to $\mathbb{N}^{2} \cup\{\infty\}$ as a partially ordered set. With respect to this identification, the semiring structure on $W_{(p)}(\mathbb{B})$ is independent of $p$, for $p>2$.

Here the connection with partitions is lost. So perhaps it is unreasonable to expect the coincidence in this case to have a deeper meaning. The semiring structure on $W_{(p)}(\mathbb{B})$ is described explicitly in (8.3) below. It is not isomorphic to either $W(\mathbb{B})$ or $W^{\text {Sch }}(\mathbb{B})$.

In Sect. 9, we establish some countability results for Witt vectors of other semirings, such as $\mathbb{N}[x]$ and the truncations $\mathbb{N} /(n=n+1)$ of $\mathbb{N}$. The final surprise of this paper is that for the truncations, there is a jump at $n=3$ :

Theorem 1.5 $W(\mathbb{N} /(n=n+1))$ is countable if and only if $n \leq 2$. The same is true for $W^{\text {Sch }}(\mathbb{N} /(n=n+1))$.

It would be interesting to improve on such cardinality statements and determine the algebraic structure of more Witt vector algebras explicitly, as we have done for Boolean Witt vectors. One could start with $\mathbb{N}[x]$ or $\mathbb{N} /(2=3)$, for example.

## 2 Background and conventions

We will follow the conventions and notations of [5]. Let us recall some of the ones that have not already appeared in the introduction.

### 2.1 Modules and algebras over $\mathbb{N}$

The set $\mathbb{N}$ of natural numbers is $\{0,1, \ldots\}$. The category $\operatorname{Mod}_{\mathbb{N}}$ of $\mathbb{N}$-modules is by definition the category of commutative monoids. The monoid operation is denoted + , and the identity is 0 . The set of $\mathbb{N}$-module homomorphisms $M \rightarrow N$ will be denoted $\operatorname{Mod}_{\mathbb{N}}(M, N)$, and similarly for other categories.

An $\mathbb{N}$-algebra (or a semiring) is an $\mathbb{N}$-module $A$ with a second commutative monoid operation $\times$ which is bilinear with respect to + , in the sense that for any $a \in A$, the map $x \mapsto a \times x$ is an endomorphism of the monoid $(A,+)$. The identity of $\times$ is
denoted 1 , and $a \times x$ is typically written $a x$. The category of $\mathbb{N}$-algebras is denoted $\mathrm{Alg}_{\mathbb{N}}$.

The Boolean semiring $\mathbb{B}$ is defined to be the quotient $\mathbb{N} /(1+1=1)$, or simply $\{0,1\}$ with $1+1=1$.

### 2.2 Partitions and symmetric functions

We will generally follow Macdonald's book [18]. In particular, a partition $\lambda$ is a sequence $\left(\lambda_{1}, \lambda_{2}, \ldots\right) \in \mathbb{N}^{\infty}$ such that $\lambda_{1} \geq \lambda_{2} \geq \cdots$ and $\lambda_{i}=0$ for sufficiently large $i$. The set of partitions will be denoted $P$. We will write $a^{b}$ for the partition $(a, a, \ldots, a, 0, \ldots)$ with $a$ occurring $b$ times. The $n$th power-sum symmetric function in $\Lambda_{\mathbb{N}}$ will be denoted by $\psi_{n}$ :

$$
\psi_{n}=\sum_{i} x_{i}^{n}
$$

instead of the more common $p_{n}$.

### 2.3 Witt vectors

As mentioned in the introduction, $W(A)$ is an $\mathbb{N}$-algebra for any $\mathbb{N}$-algebra $A$. The binary operations,$+ \times$ on $W(A)$ are induced by co-operations on the representing object $\Lambda_{\mathbb{N}}$ :

$$
\begin{equation*}
\Delta^{+}, \Delta^{\times}: \Lambda_{\mathbb{N}} \longrightarrow \Lambda_{\mathbb{N}} \otimes_{\mathbb{N}} \Lambda_{\mathbb{N}} \tag{2.3.1}
\end{equation*}
$$

In terms of symmetric functions in the formal variables $x_{1}, x_{2}, \ldots$, they are given by the following rules:

$$
\begin{aligned}
& \Delta^{+}: f\left(\ldots, x_{i}, \ldots\right) \mapsto f\left(\ldots, x_{i} \otimes 1,1 \otimes x_{i}, \ldots\right) \\
& \Delta^{\times}: f\left(\ldots, x_{i}, \ldots\right) \mapsto f\left(\ldots, x_{i} \otimes x_{j}, \ldots\right)
\end{aligned}
$$

For example, we have

$$
\begin{aligned}
& \Delta^{+}\left(\psi_{n}\right)=\psi_{n} \otimes 1+1 \otimes \psi_{n} \\
& \Delta^{\times}\left(\psi_{n}\right)=\psi_{n} \otimes \psi_{n}
\end{aligned}
$$

These equations are the properties of $\Delta^{+}$and $\Delta^{\times}$we will most often use, and in fact they determine $\Delta^{+}$and $\Delta^{\times}$uniquely and hence furnish alternative definitions.

Similarly, the additive and multiplicative units $0,1 \in W(A)$ are induced by co-units

$$
\begin{equation*}
\varepsilon^{+}, \varepsilon^{\times}: \Lambda_{\mathbb{N}} \longrightarrow \mathbb{N} \tag{2.3.2}
\end{equation*}
$$

given by

$$
\begin{aligned}
& \varepsilon^{+}: f\left(\ldots, x_{i}, \ldots\right) \mapsto f(0,0,0, \ldots) \\
& \varepsilon^{\times}: f\left(\ldots, x_{i}, \ldots\right) \mapsto f(1,0,0, \ldots) .
\end{aligned}
$$

The same statements hold for $\Lambda^{\text {Sch }}$ instead of $\Lambda_{\mathbb{N}}$.
For any $x \in A$, the anti-Teichmüller lift $\{x\} \in W^{\operatorname{Sch}}(A)$ is defined by

$$
\{x\}: s_{\lambda} \mapsto \begin{cases}x^{r} & \text { if } \lambda=1^{r}  \tag{2.3.3}\\ 0 & \text { otherwise } .\end{cases}
$$

So for example, we have

$$
\begin{equation*}
\sigma(\{x\})=1+x t+x^{2} t^{2}+\cdots=(1-x t)^{-1} . \tag{2.3.4}
\end{equation*}
$$

See [5, (6.10)] for the basic properties of anti-Teichmüller lifts. (Note that the symbol $\sigma$ denotes a different map in [5], corresponding to different sign conventions. There two maps $\sigma_{+}^{+}$and $\sigma_{-}^{-}$are defined. Here $\sigma$ means $\sigma_{+}^{+}$, whereas there it means $\sigma_{-}^{-}$.)

We end with some words on composition structures, which will have a small part in this paper. Consider a composition $\mathbb{N}$-algebra $Q$, such as $\Lambda_{\mathbb{N}}$ or $\Lambda^{\text {Sch }}$. (See (3.3) in [5] for the general definition.) Then the $Q$-Witt vectors $W_{Q}(A)=\operatorname{Alg}_{\mathbb{N}}(Q, A)$ have not only an algebra structure but also an action of $Q$. For $f \in Q, a \in W_{Q}(A)$, the element $f(a) \in W_{Q}(A)$ is the map $Q \rightarrow A$ given by

$$
f(a): g \mapsto a(g \circ f),
$$

where o denotes the composition operation on $Q$-this is simply plethysm of symmetric functions in all examples in this paper. In fact, $W_{Q}$ is the right adjoint of the forgetful functor from $Q$-semirings to semirings. This generalizes the fact that the big Witt functor is the right adjoint of the forgetful functor from $\lambda$-rings to rings.

## 3 Orders on Witt vectors

This section introduces some generalities on preorders. They have only a small part in this paper; so it can be skipped and consulted as needed.

## 3.1 $\mathrm{Alg}_{\mathbb{N}}-$ preorders

Recall that a preorder is a reflexive transitive relation. Let us say that a relation $\leq$ on an $\mathbb{N}$-algebra $A$ is an $\mathrm{Alg}_{\mathbb{N}}$-preorder if it is a preorder and satisfies the implication

$$
\begin{equation*}
a \leq b \text { and } a^{\prime} \leq b^{\prime} \Longrightarrow a+a^{\prime} \leq b+b^{\prime} \text { and } a a^{\prime} \leq b b^{\prime} \tag{3.1.1}
\end{equation*}
$$

This implication is equivalent to requiring that the graph of $\leq$ be a sub- $\mathbb{N}$-algebra of $A \times A$. Hence an $\mathrm{Alg}_{\mathbb{N}}$-preorder is the same as a preorder object over $A$ in the category $\mathrm{Alg}_{\mathbb{N}}$.

The most important example for this paper is the difference preorder defined by

$$
\begin{equation*}
a \leq b \Longleftrightarrow \exists c \in A \text { such that } b=a+c . \tag{3.1.2}
\end{equation*}
$$

It is the strongest $\mathrm{Alg}_{\mathbb{N}}$-preorder such that $c \geq 0$ for all $c \in A$. A second example is $a \leq b \Leftrightarrow a=b$; a third is $a \leq b$ for all $a, b$, which agrees with (3.1.2) when $A$ is a ring.

### 3.2 Preorders on Witt vectors

Let $Q$ be a composition $\mathbb{N}$-algebra, for example $\Lambda_{\mathbb{N}}$ or $\Lambda^{\text {Sch }}$. (See (3.3) in [5] for the general definition.) Then given any preorder $\leq$ on an $\mathbb{N}$-algebra $A$, we can define a preorder on the algebra $W_{Q}(A)=\operatorname{Hom}_{\operatorname{Alg}_{\mathbb{N}}}(Q, A)$ of Witt vectors: For $a, b \in W_{Q}(A)$, we write

$$
\begin{equation*}
a \preccurlyeq b \Longleftrightarrow a(x) \leq b(x) \quad \text { for all } x \in Q . \tag{3.2.1}
\end{equation*}
$$

Observe that if $\leq$ is a partial order (that is, it satisfies $a \leq b \leq a \Rightarrow a=b$ ), then so is $\preccurlyeq$.

Now assume that $\leq$ is an Alg $_{\mathbb{N}^{-}}$preorder. Then $\preccurlyeq$ is also an Alg $_{\mathbb{N}}$-preorder. Indeed, if $R \subseteq A \times A$ denotes the graph of $\leq$, then $W_{Q}(R)$ is a subalgebra of $W_{Q}(A \times A)=$ $W_{Q}(A) \times W_{Q}(A)$ and is clearly the graph of $\preccurlyeq$. In fact, $W_{Q}(R)$ is a sub- $Q$-semiring of $W_{Q}(A) \times W_{Q}(A)$. So $\preccurlyeq$ is a preorder in the category of $Q$-semirings. Equivalently, it is an Alg $_{\mathbb{N}}$-preorder satisfying

$$
\begin{equation*}
a \preccurlyeq b \Longrightarrow f(a) \preccurlyeq f(b) \tag{3.2.2}
\end{equation*}
$$

for any $f \in Q$.
Finally observe that if a subset $S \subseteq Q$ generates $Q$ as an $\mathbb{N}$-algebra, then for all $a, b \in W_{Q}(A)$, we have

$$
\begin{equation*}
a \preccurlyeq b \Longleftrightarrow a(s) \leq b(s) \quad \text { for all } s \in S \tag{3.2.3}
\end{equation*}
$$

Indeed, since $R$ is a subalgebra of $A \times A$, the map $(a, b): Q \rightarrow A \times A$ factors through $R$ if and only if the image of $S$ is contained in $R$.

### 3.3 The preorders in this paper

The preorder $\leq$ on $\mathbb{N}$-algebras will always be understood to be the difference preorder of (3.1.2), and the preorder $\preccurlyeq$ on Witt vectors $W_{Q}(A)$ will always be understood to be the one induced by the difference preorder on $A$.

The following result implies that for all the composition $\mathbb{N}$-algebras $Q$ that appear in this paper-namely $\Lambda_{\mathbb{N}}, \Lambda^{\text {Sch }}$, and $\Lambda_{\mathbb{N},(p)}$-the difference preorder $\leq$ on $W_{Q}(A)$ is stronger than $\preccurlyeq$.

Proposition 3.4 For $a, c \in W_{Q}(A)$, we have

$$
\begin{equation*}
a \preccurlyeq a+c, \tag{3.4.1}
\end{equation*}
$$

as long as $Q$ has the property that for all $x \in Q$, we have $\varepsilon^{+}(x) \leq x$, where $\varepsilon^{+}: Q \rightarrow$ $\mathbb{N}$ denotes the additive co-unit of $Q$.
Proof Given any element $z=\sum_{i} x_{i} \otimes y_{i} \in Q \otimes_{\mathbb{N}} Q$, write $x_{i}=\varepsilon^{+}\left(x_{i}\right)+x_{i}^{\prime}$, for some elements $x_{i}^{\prime} \in Q$. Then we have

$$
z=\sum_{i}\left(\varepsilon^{+}\left(x_{i}\right)+x_{i}^{\prime}\right) \otimes y_{i}=\sum_{i} \varepsilon^{+}\left(x_{i}\right) \otimes y_{i}+\sum_{i} x_{i}^{\prime} \otimes y_{i} \geq\left(\varepsilon^{+} \otimes \mathrm{id}\right)(z)
$$

When $z=\Delta^{+}(x)$ for some $x \in Q$, this simplifies to $\Delta^{+}(x) \geq 1 \otimes x$. Applying the homomorphism $c \otimes a: Q \otimes_{\mathbb{N}} Q \rightarrow A$ to this inequality, we obtain $(c+a)(x) \geq a(x)$. Since this holds for all $x \in Q$, we have $c+a \succcurlyeq a$.

## $4 W(\mathbb{B})$

The purpose of this section is to determine $W(\mathbb{B})$.

### 4.1 Partition ideals

Recall that $P$ denotes the set of partitions. Let us say that a subset $I \subseteq P$ is a partition ideal if the following implication holds for all $\lambda, \mu \in P$ :

$$
\lambda \in I, \quad \lambda \subseteq \mu \Longrightarrow \mu \in I .
$$

Here, $\subseteq$ stands for the inclusion ordering of partitions, i.e., two partitions $\lambda$ and $\mu$ satisfy $\lambda \subseteq \mu$ if and only if $\lambda_{i} \leq \mu_{i}$ for all $i$ (including those for which $\mu_{i}=0$ ).

For any symmetric function $g \in \Lambda_{\mathbb{N}}$, let us say the (monomial) constituents of $g$ are the partitions $v$ such that $g_{\nu} \neq 0$ in the unique decomposition $g=\sum_{\nu} g_{\nu} m_{\nu}$, $g_{v} \in \mathbb{N}$, where the $m_{v}$ are the monomial symmetric functions. The reason for the optional modifier monomial here and below is that we will consider Schur analogues of these concepts in the next section.

Let us say that a partition ideal $I$ is (monomial) prime if $I \neq P$ and whenever partitions $\lambda, \mu \in P$ have the property that every constituent of $m_{\lambda} m_{\mu}$ is in $I$, then either $\lambda$ or $\mu$ lies in $I$.

For any linear map $f \in \operatorname{Mod}_{\mathbb{N}}\left(\Lambda_{\mathbb{N}}, \mathbb{B}\right)$, write

$$
\operatorname{pker}(f)=\left\{\lambda \in P \mid f\left(m_{\lambda}\right)=0\right\}
$$

We will call $\operatorname{pker}(f)$ the (monomial) partition kernel of $f$.
Understanding $W(\mathbb{B})$ hinges on knowing the constituents of the product of two monomial symmetric functions. This is a combinatorial problem which is not very deep:

Proposition 4.2 Let $\lambda$ and $\mu$ be two partitions. Then the constituents of the product $m_{\lambda} m_{\mu}$ are the partitions $v$ which can be obtained by picking a permutation $\sigma$ of the set $\{1,2,3, \ldots\}$ and rearranging the sequence

$$
\left(\lambda_{1}+\mu_{\sigma(1)}, \lambda_{2}+\mu_{\sigma(2)}, \lambda_{3}+\mu_{\sigma(3)}, \ldots\right)
$$

in weakly decreasing order.
Proof Clear.
Corollary 4.3 Let $\lambda$ and $\mu$ be two partitions. Then, every constituent $v$ of $m_{\lambda} m_{\mu}$ satisfies $\lambda \subseteq \nu$ and $\mu \subseteq \nu$.

Proof Let $\nu$ be a constituent of $m_{\lambda} m_{\mu}$. Due to (4.2), there exists a permutation $\sigma$ of the set $\{1,2,3, \ldots\}$ such that $v$ is the weakly decreasing permutation of the sequence $\left(\lambda_{1}+\mu_{\sigma(1)}, \lambda_{2}+\mu_{\sigma(2)}, \lambda_{3}+\mu_{\sigma(3)}, \ldots\right)$. But for each positive integer $i$, at least $i$ terms of the sequence $\left(\lambda_{1}+\mu_{\sigma(1)}, \lambda_{2}+\mu_{\sigma(2)}, \lambda_{3}+\mu_{\sigma(3)}, \ldots\right)$ are $\geq \lambda_{i}$ (since at least $i$ terms of the sequence $\left(\lambda_{1}, \lambda_{2}, \lambda_{3}, \ldots\right)=\lambda$ are $\left.\geq \lambda_{i}\right)$. Hence, for each positive integer $i$, at least $i$ terms of $v$ (which is a permutation of this sequence) are $\geq \lambda_{i}$. Since $\nu$ is a partition, this entails $\nu_{i} \geq \lambda_{i}$. Thus, $\lambda \subseteq v$, and similarly $\mu \subseteq \nu$.

The next proposition is more technical and will be used to classify the prime partition ideals:

Proposition 4.4 Let $\lambda$ and $\mu$ be two partitions, and $i$ and $j$ be two positive integers such that $\lambda_{i}>\mu_{i}, \lambda_{i}>\lambda_{i+1}, \mu_{j}>\lambda_{j}$, and $\mu_{j}>\mu_{j+1}$. Let $\bar{\lambda}$ denote the result of diminishing the ith part of $\lambda$ by 1 , and let $\bar{\mu}$ be the result of diminishing the $j$ th part of $\mu$ by 1 . Then, for every constituent $\nu$ of $m_{\bar{\lambda}} m_{\bar{\mu}}$, we have either $\lambda \subseteq \nu$ or $\mu \subseteq \nu$.

Proof Since $\lambda_{i}>\mu_{i}$ and $\mu_{j}>\lambda_{j}$, we have $i \neq j$. Thus, we can assume by symmetry that $i<j$.

Fix a constituent $v$ of $m_{\bar{\lambda}} m_{\bar{\mu}}$. Putting (4.2) into words, we see that the constituents of $m_{\bar{\lambda}} m_{\bar{\mu}}$ are found by combining rows of $\bar{\lambda}$ with rows (possibly with zero boxes) of $\bar{\mu}$, such that no two rows of $\bar{\lambda}$ are combined with the same row of $\bar{\mu}$ and vice versa. Let us consider the matching between rows of $\bar{\lambda}$ and rows of $\bar{\mu}$ that gives rise to the constituent $v$. We distinguish between two cases, depending on whether (i) the $i$ th row of $\bar{\lambda}$ is combined with a nonzero-length row of $\bar{\mu}$ or (ii) with a zero-length row of $\bar{\mu}$.

Case ( $i$ ). If the $i$ th row of $\bar{\lambda}$ is combined with a row of $\bar{\mu}$ with at least one box, then the resulting row has at least $\bar{\lambda}_{i}+1=\lambda_{i}$ boxes. Hence, in this case, the partition $v$ has at least $i$ rows of length $\geq \lambda_{i}$ (namely, the row obtained by combining the $i$ th row of $\bar{\lambda}$ with its match, and also the rows obtained by combining the first $i-1$ rows of $\bar{\lambda}$ with their matches). In other words, $v_{i} \geq \lambda_{i}$. Together with $\bar{\lambda} \subseteq v$ (which follows from (4.3)), this yields $\lambda \subseteq \nu$, which completes the proof in case (i).

Case (ii). If the $i$ th row of $\bar{\lambda}$ is combined with a row of $\bar{\mu}$ with zero boxes, then the resulting row has $\bar{\lambda}_{i} \geq \mu_{i} \geq \mu_{j}$ boxes. But the first $j-1$ rows of $\bar{\mu}$, whatever they are combined with, also have at least $\bar{\mu}_{j-1}=\mu_{j-1} \geq \mu_{j}$ boxes. So there are at least
$j$ rows in $\nu$ with at least $\mu_{j}$ boxes. Consequently, $\nu_{j} \geq \mu_{j}$. Together with $\bar{\mu} \subseteq v$ (this is a consequence of (4.3)), this results in $\mu \subseteq \nu$, which completes the proof in case (ii).

Next we include a fact we will not use until Sect. 9:
Proposition 4.5 Let $\lambda$ and $\mu$ be two partitions.
(1) Let $\lambda+\mu$ denote the partition $\left(\lambda_{1}+\mu_{1}, \lambda_{2}+\mu_{2}, \lambda_{3}+\mu_{3}, \ldots\right)$. Then, $\lambda+\mu$ is a constituent of $m_{\lambda} m_{\mu}$.
(2) Let $\lambda \uplus \mu$ denote the partition obtained by sorting the components of the vector $\left(\lambda_{1}, \mu_{1}, \lambda_{2}, \mu_{2}, \lambda_{3}, \mu_{3}, \ldots\right)$ in weakly decreasing order. Then, $\lambda \uplus \mu$ is a constituent of $m_{\lambda} m_{\mu}$.
(3) If $\lambda$ and $\mu$ are both nonzero, then $m_{\lambda} m_{\mu}$ has at least two distinct nonzero constituents.

Proof (1) Apply (4.2) where $\sigma$ is the identity permutation.
(2) Apply (4.2) where $\sigma$ has the property that $\sigma(1), \ldots, \sigma(b)$ are all greater than the length of $\lambda$, where $b$ is the length of $\mu$.
(3) Since $\lambda$ and $\mu$ are both nonzero, the partitions $\lambda+\mu$ and $\lambda \uplus \mu$ have different lengths, and are thus distinct. (The length of the former is the maximum of the lengths of $\lambda$ and $\mu$, while the length of the latter is the sum.) They are also nonzero and, by parts (1) and (2), constituents of $m_{\lambda} m_{\mu}$.
Proposition 4.6 Let $2^{P}$ denote the set of subsets of $P$.
(1) The map $f \mapsto \operatorname{pker}(f)$ is a bijection $\operatorname{Mod}_{\mathbb{N}}\left(\Lambda_{\mathbb{N}}, \mathbb{B}\right) \rightarrow 2^{P}$.
(2) $\operatorname{pker}(f)$ is a partition ideal if and only if $f$ satisfies $f(x)=0 \Rightarrow f(x y)=0$.
(3) $\operatorname{pker}(f)$ is a prime partition ideal if and only if $f$ is an $\mathbb{N}$-algebra map.

Proof (1) Linear maps $\Lambda_{\mathbb{N}} \rightarrow \mathbb{B}$ are in bijection with set maps $P \rightarrow \mathbb{B}$, which are in bijection with subsets of $P$.
(2) From (4.2), it follows that for any partition $\lambda$, we have

$$
m_{1} m_{\lambda}=\sum_{\mu} a_{\mu} m_{\mu}
$$

where the sum is over all partitions $\mu$ obtained by adding a single box to (the Young diagram of) $\lambda$, and the $a_{\mu}$ are at least 1 . Therefore we have

$$
f\left(m_{1} m_{\lambda}\right)=\sum_{\mu} a_{\mu} f\left(m_{\mu}\right) .
$$

If $\lambda \in \operatorname{pker}(f)$ and if $f$ satisfies $f(x)=0 \Rightarrow f(x y)=0$, then we have $f\left(m_{1} m_{\lambda}\right)=0$ and hence $f\left(m_{\mu}\right)=0$ for all $\mu$ obtained by adding one box to $\lambda$. It then follows by induction that $f\left(m_{\mu}\right)=0$ for all $\mu$ with $\lambda \subseteq \mu$, which means that $\operatorname{pker}(f)$ is a partition ideal.

Conversely, suppose $\operatorname{pker}(f)$ is a partition ideal. Let $\lambda \in \operatorname{pker}(f)$. Since for any partition $\mu$, every constituent $\nu$ of $m_{\lambda} m_{\mu}$ satisfies $\lambda \subseteq \nu$ (by (4.3)), we have
$f\left(m_{\lambda} m_{\mu}\right)=0$. The general implication $f(x)=0 \Rightarrow f(x y)=0$ then follows. Indeed, write $x=\sum_{\lambda} x_{\lambda} m_{\lambda}, y=\sum_{\mu} y_{\mu} m_{\mu}$. Since $f(x)=0$, it follows that whenever $x_{\lambda} \neq 0$, we have $\lambda \in \operatorname{pker}(f)$ and hence $f\left(m_{\lambda} m_{\mu}\right)=0$. Therefore we have

$$
f(x y)=\sum_{\lambda, \mu} x_{\lambda} y_{\mu} f\left(m_{\lambda} m_{\mu}\right)=0
$$

(3) First observe that under either assumption, we have $\operatorname{pker}(f) \neq P$ and $f(1)=1$. Indeed, if $f$ is an $\mathbb{N}$-algebra map, then we have $1=f(1)=f\left(m_{0}\right)$ and hence $\operatorname{pker}(f) \neq P$. Conversely, if $\operatorname{pker}(f)$ is a prime partition ideal, then there is a partition $\lambda$ such that $\lambda \notin \operatorname{pker}(f)$. Since $\operatorname{pker}(f)$ is an ideal, we have $0 \notin \operatorname{pker}(f)$ and hence that $f(1)=f\left(m_{0}\right)=1$.

Now suppose that $f$ is an $\mathbb{N}$-algebra map and that $\lambda, \mu$ are partitions such that every constituent of $m_{\lambda} m_{\mu}$ is in $\operatorname{pker}(f)$. Then we have $f\left(m_{\lambda}\right) f\left(m_{\mu}\right)=f\left(m_{\lambda} m_{\mu}\right)=0$. It follows that either $\lambda \in \operatorname{pker}(f)$ or $\mu \in \operatorname{pker}(f)$, and hence that $\operatorname{pker}(f)$ is prime.

Conversely, suppose $\operatorname{pker}(f)$ is prime. By part (2), it is enough to show that if $f(x)=f(y)=1$ then $f(x y)=1$. In this case, $x$ has a constituent $\lambda \notin \operatorname{pker}(f)$ and $y$ has a constituent $\mu \notin \operatorname{pker}(f)$. Since $\operatorname{pker}(f)$ is prime, $m_{\lambda} m_{\mu}$ must have a constituent not in $\operatorname{pker}(f)$. This is also a constituent of $x y$. Therefore $f(x y)=1$.

Proposition 4.7 For integers $x, y \in \mathbb{N}$, consider the partition ideal

$$
I_{(x, y)}=\left\{\lambda \in P \mid(y+1)^{x+1} \subseteq \lambda\right\},
$$

where $(y+1)^{x+1}$ denotes the partition $(y+1, \ldots, y+1)$ with $y+1$ occurring $x+1$ times. Then $I_{(x, y)}$ is prime. Conversely, every nonempty prime partition ideal is of this form.

Proof Let us first show that $I_{(x, y)}$ is prime. Let $\lambda$ and $\mu$ be partitions, and suppose all the constituents of $m_{\lambda} m_{\mu}$ are in $I_{(x, y)}$. By symmetry, we may assume $\lambda_{x+1} \geq \mu_{x+1}$. We will show $\lambda \in I_{(x, y)}$. Let $v$ denote the partition obtained by sorting the components of the vector

$$
\left(\lambda_{1}+\mu_{1}, \ldots, \lambda_{x}+\mu_{x}, \lambda_{x+1}, \mu_{x+1}, \lambda_{x+2}, \mu_{x+2}, \ldots\right)
$$

in weakly decreasing order. Then $v$ is a constituent of $m_{\lambda} m_{\mu}$ (due to (4.2)), and hence we have $v \in I_{(x, y)}$. Further the sorting does not change the first $x+1$ components, and so we have $\lambda_{x+1}=v_{x+1} \geq y+1$ and hence $\lambda \in I_{(x, y)}$.

For the converse, let $I$ be a nonempty prime partition ideal. Let us first show that $I$ contains a unique minimal element with respect to $\subseteq$. Since $I$ is nonempty, there exists some minimal element. Now suppose for a contradiction that we have two distinct minimal elements $\lambda, \mu \in I$. Then $\lambda \nsubseteq \mu$, and so there exists an integer $i$ such that $\lambda_{i}>\mu_{i}$ and $\lambda_{i}>\lambda_{i+1}$. Likewise, there is a $j$ such that $\mu_{j}>\lambda_{j}$ and $\mu_{j}>\mu_{j+1}$. Now let $\bar{\lambda}$ denote the result of diminishing the $i$ th part of $\lambda$ by 1 , and let $\bar{\mu}$ be the result of diminishing the $j$ th part of $\mu$ by 1 . By (4.4), for every constituent $v$ of $m_{\bar{\lambda}} m_{\bar{\mu}}$, we have either $\lambda \subseteq v$ or $\mu \subseteq v$, and hence $v \in I$. Therefore, since the partition ideal $I$ is prime, either $\bar{\lambda}$ or $\bar{\mu}$ must lie in $I$, and so either $\lambda$ or $\mu$ is not minimal.

It remains to show that the minimal element $\pi$ of $I$ is of the form $(y+1)^{x+1}$. This is equivalent to showing that if $\pi$ is contained in the union of two partitions $\lambda$ and $\mu$ (that is, we have $\pi_{i} \leq \lambda_{i}$ or $\pi_{i} \leq \mu_{i}$ for all $i$ ), then we have $\pi \subseteq \lambda$ or $\pi \subseteq \mu$. But for any constituent $\nu$ of $m_{\lambda} m_{\mu}$, we have $\lambda \subseteq \nu$ and $\mu \subseteq \nu$, by (4.3). It follows that $\pi_{i} \leq v_{i}$ for all $i$. In other words, $v$ lies in $I$. Since $I$ is prime, we have $\lambda \in I$ or $\mu \in I$, and hence $\pi \subseteq \lambda$ or $\pi \subseteq \mu$.

### 4.8 Explicit description of $W(\mathbb{B})$

Here we prove parts (1)-(2) of Theorem (1.1). Following the notation there, for $x, y \in$ $\mathbb{N}$, let $(x, y)$ denote the element of $W(\mathbb{B})=\operatorname{Alg}_{\mathbb{N}}\left(\Lambda_{\mathbb{N}}, \mathbb{B}\right)$ corresponding under (4.6) to the prime partition ideal $I_{(x, y)}$. Thus we have

$$
(x, y): m_{\lambda} \mapsto \begin{cases}0 & \text { if } \lambda_{x+1} \geq y+1 \\ 1 & \text { otherwise }\end{cases}
$$

Also let $\infty \in W(\mathbb{B})$ denote the element corresponding to the empty partition ideal, which is vacuously prime; so $\infty: m_{\lambda} \mapsto 1$. This defines a map

$$
\begin{equation*}
\mathbb{N}^{2} \cup\{\infty\} \xrightarrow{\sim} W(\mathbb{B}) \tag{4.8.1}
\end{equation*}
$$

which is a bijection by (4.6)-(4.7). In these terms, the $\Lambda_{\mathbb{N}}$-semiring structure on $W$ ( $\mathbb{B}$ ) is given by the following proposition.

Proposition 4.9 (1) Under the unique $\mathbb{N}$-algebra map $\mathbb{N} \rightarrow W(\mathbb{B})$, the image of an element $x$ is $(x, 0)$. In particular, we have $0=(0,0), 1=(1,0)$, and $(x, 0)+$ $\left(x^{\prime}, 0\right)=\left(x+x^{\prime}, 0\right)$ for all $x, x^{\prime} \in \mathbb{N}$.
(2) We have $(x, 0)+(0, y)=(x, y)$, for all $x, y \in \mathbb{N}$.
(3) We have $(0, y)+\left(0, y^{\prime}\right)=\left(0, \max \left\{y, y^{\prime}\right\}\right)$ and $(0, y)\left(0, y^{\prime}\right)=\left(0, \min \left\{y, y^{\prime}\right\}\right)$, for all $y, y^{\prime} \in \mathbb{N}$.
(4) We have $\infty+z=\infty$ for all $z \in W$ ( $\mathbb{B})$.
(5) We have $\infty \cdot \infty=\infty$ and $\infty \cdot(0, y)=(0, y)$ for all $y \in \mathbb{N}$.
(6) For any nonzero partition $\lambda$, we have $m_{\lambda}(\infty)=\infty$ and $m_{\lambda}(0, y)=\left(0,\left[y / \lambda_{1}\right]\right)$, where $\left[y / \lambda_{1}\right]$ denotes the greatest integer at most $y / \lambda_{1}$.
(7) The bijection (4.8.1) identifies the partial order $\preccurlyeq$ on $W(\mathbb{B})$ with the componentwise partial order on $\mathbb{N}^{2} \cup\{\infty\}$. That is, we have $\left(x_{1}, y_{1}\right) \preccurlyeq\left(x_{2}, y_{2}\right)$ if and only if $x_{1} \leq x_{2}$ and $y_{1} \leq y_{2}$, and $z \preccurlyeq \infty$ for all $z \in W(\mathbb{B})$.

Proof (1) Let $f$ denote the image of $x$ under the map $\mathbb{N} \rightarrow W(\mathbb{N})$. Then we have $\sigma(f)=\sigma(1)^{x}=(1+t)^{x}=1+\cdots+t^{x}$. Thus $f: \Lambda_{\mathbb{N}} \rightarrow \mathbb{N}$ satisfies $f\left(m_{1^{x+1}}\right)=$ $f\left(e_{x+1}\right)=0$ and $f\left(m_{1^{x}}\right)=f\left(e_{x}\right)=1$. Therefore its image in $W(\mathbb{B})$ has the same properties; so its image is $(x, 0)$.
(2) Let $I$ denote the partition kernel of the homomorphism $\Lambda_{\mathbb{N}} \rightarrow \mathbb{B}$ given by $(x, 0)+(0, y)$. By the above, $I$ is empty or is of the form $I_{\left(x^{\prime}, y^{\prime}\right)}$. So it is enough to show that the partition $(y+1)^{x+1}$ is an element of $I$ but that $y^{x+1}$ and $(y+1)^{x}$ are not.

Recall that we have

$$
\Delta^{+}\left(m_{\lambda}\right)=\sum_{\mu, v} m_{\mu} \otimes m_{v}
$$

where $\mu$ and $v$ run over all partitions such that the vector $\left(\mu_{1}, \nu_{1}, \mu_{2}, \nu_{2}, \ldots\right)$ is a permutation of $\lambda$. Therefore $\lambda$ is in the kernel of $(x, 0)+(0, y)$ if and only if for all such decompositions of $\lambda$, we have either $1^{x+1} \subseteq \mu$ or $(y+1)^{1} \subseteq \nu$.

Applying this to $\lambda=(y+1)^{x+1}$, we get terms of the form $m_{(y+1)^{i}} \otimes m_{(y+1)^{j}}$ where $i+j=x+1$. If $i>0$, then $(y+1)^{i}$ contains $(y+1)^{1}$. Otherwise, $j=x+1$, in which case $(y+1)^{x+1}$ contains $1^{x+1}$. So $(y+1)^{x+1}$ is in the partition kernel $I$.

On the other hand, we can decompose $\lambda=y^{x+1}$ into $\mu=0$ and $\nu=y^{x+1}$. Then $\mu$ does not contain $1^{x+1}$ and $v$ does not contain $(y+1)^{1}$. So $y^{x+1}$ is not in $I$. Similarly $(y+1)^{x}$ decomposes into $\mu=(y+1)^{x}$ and $v=0$, the first of which does not contain $1^{x+1}$ and the second of which does not contain $(y+1)^{1}$. So $(y+1)^{x}$ is not in $I$.
(3) The image of any power sum $\psi_{r}=\sum_{i} x_{i}^{r}$ under $\Delta^{+}$is $\psi_{r} \otimes 1+1 \otimes \psi_{r}$. Therefore the partition $r^{1}$ is in the partition kernel of $(0, y)+\left(0, y^{\prime}\right)$ if and only if $r$ is greater than $y$ and $y^{\prime}$. Therefore we have $(0, y)+\left(0, y^{\prime}\right)=\left(0, \max \left\{y, y^{\prime}\right\}\right)$.

Similarly, we have $\Delta^{\times}\left(\psi_{r}\right)=\psi_{r} \otimes \psi_{r}$, and so the partition $r^{1}$ is in the partition kernel of $(0, y)\left(0, y^{\prime}\right)$ if and only if $r$ is greater than $y$ or $y^{\prime}$. Therefore we have $(0, y)\left(0, y^{\prime}\right)=\left(0, \min \left\{y, y^{\prime}\right\}\right)$.
(4) Since $\infty+z \geq \infty$, we have $\infty+z \succcurlyeq \infty$, by (3.4), and hence $\infty+z=\infty$.
(5) The map $\Delta^{\times}$is injective because it has a retraction, for example the multiplicative co-unit $\varepsilon^{\times}$. Therefore for any partition $\lambda$, the tensor $\Delta^{\times}\left(m_{\lambda}\right)$ is a nonempty sum of basic tensors $m_{\mu} \otimes m_{\nu}$, all of which pair with $(\infty, \infty)$ to make 1 . Therefore we have $(\infty \cdot \infty)\left(m_{\lambda}\right)=1$ for all $\lambda$, and hence $\infty \cdot \infty=\infty$.

We have $\Delta^{\times}\left(\psi_{r}\right)=\psi_{r} \otimes \psi_{r}$ which pairs with $((0, y), \infty)$ to make $(0, y)\left(\psi_{r}\right)$, which is 0 if and only if $r>y$. Therefore the partition $r^{1}$ is in $\operatorname{pker}(\infty \cdot(0, y))$ if and only if $r>y$. Thus we have $\operatorname{pker}(\infty \cdot(0, y))=I_{(0, y)}$, and hence $(0, y) \cdot \infty=(0, y)$.
(6) The value of $m_{\lambda}(\infty) \in W(\mathbb{B})=\operatorname{Alg}_{\mathbb{N}}\left(\Lambda_{\mathbb{N}}, \mathbb{B}\right)$ at $m_{\mu} \in \Lambda_{\mathbb{N}}$ is by definition the value of $\infty$ at $m_{\mu} \circ m_{\lambda}$. Since $\lambda$ is nonzero, $m_{\lambda}$ is a sum of infinitely many monomials in $x_{1}, x_{2}, \ldots$ Therefore $m_{\mu} \circ m_{\lambda}$ is the sum of at least one such monomial and is hence nonzero. It follows that $\infty\left(m_{\mu} \circ m_{\lambda}\right)=1$ for all $\mu$ and hence $m_{\lambda}(\infty)=\infty$.

Similarly, the value of $m_{\lambda}(0, y)$ at $\psi_{r}$ is $(0, y)\left(\psi_{r} \circ m_{\lambda}\right)$. But this is 1 if and only if $r \lambda_{1} \leq y$. Therefore $\left(m_{\lambda}(0, y)\right)\left(\psi_{r}\right)$ is 1 if and only if $r \leq\left[y / \lambda_{1}\right]$. And so we have $m_{\lambda}(0, y)=\left(0,\left[y / \lambda_{1}\right]\right)$.
(7) Both sides of the first part are clearly equivalent to $I_{\left(x_{1}, y_{1}\right)} \supseteq I_{\left(x_{2}, y_{2}\right)}$. The second part holds because $\infty$ corresponds to the empty partition ideal.

### 4.10 An interpretation of the algebraic structure on $W(\mathbb{B})$

We can express the description of $W(\mathbb{B})$ given above in terms of the Dorroh extension, a general construction in commutative algebra over $\mathbb{N}$. Let us first recall it.

Let $f: A \rightarrow B$ be a homomorphism of $\mathbb{N}$-algebras. The Dorroh extension $D(f)$ is defined to be $A \times B$ with the usual, product $\mathbb{N}$-module structure and with a multiplication law given by the formula

$$
(x, y)\left(x^{\prime}, y^{\prime}\right)=\left(x x^{\prime}, f(x) y^{\prime}+y f\left(x^{\prime}\right)+y y^{\prime}\right) .
$$

One can check that this is an $\mathbb{N}$-algebra structure on $D(f)$ with multiplicative identity $(1,0)$. For example, if $B$ is additively cancellative, then the map $D(f) \rightarrow A \times B$ sending $(x, y)$ to $(x, f(x)+y)$ is an injective homomorphism with image $\{(a, b) \mid f(a) \leq$ $b\}$.

Now let $\mathbb{N}^{\prime}$ denote $\mathbb{N} \cup\{\infty\}$ with the $\mathbb{N}$-algebra structure where addition is max and multiplication is min. Let $f$ denote the unique $\mathbb{N}$-algebra map $\mathbb{N} \rightarrow \mathbb{N}^{\prime}$. (So $f(n)=\infty$ unless $n=0$.) Define a set map $g: D(f) \rightarrow W(\mathbb{B})$ by $(x, y) \mapsto(x, y)$ if $y \neq \infty$ and $(x, \infty) \mapsto \infty$. Then $g$ is a surjective $\mathbb{N}$-algebra morphism. It follows that $W(\mathbb{B})$ is the quotient of $D(f)$ obtained by identifying all elements of the form $(x, \infty)$.

The algebra $\mathbb{N}^{\prime}$ can perhaps be demystified by observing that it is isomorphic to the image of the ghost map $w: W(\mathbb{B}) \rightarrow \mathbb{B}^{\infty}$. (See (6.2.1) of [5].) Indeed, for a Witt vector of the form $(0, y)$, its image under the ghost map is the bit vector $\left\langle z_{1}, z_{2}, \ldots\right\rangle$ with $z_{r}=1$ if $r \leq y$ and otherwise $z_{r}=0$. Every other Witt vector has image $\langle 1,1, \ldots\rangle$. Clearly addition and multiplication of such bit vectors are given by taking the maximum and minimum of the corresponding $y$ values. One might say that $W(\mathbb{B})$ can in essence be recovered by applying a general construction to the image of the ghost map. It would be interesting if something like this holds more generally.

It would also be interesting to give an explicit description of $W(A)$ as a $\Lambda_{\mathbb{N}}$-semiring for any $\mathbb{B}$-algebra $A$. This would be tantamount to having an explicit understanding of multiplication, the coproducts $\Delta^{+}$and $\Delta^{\times}$, and all right plethysm maps $f \mapsto f \circ m_{\lambda}$ on the semiring $\mathbb{B} \otimes_{\mathbb{N}} \Lambda_{\mathbb{N}}=\bigoplus_{\lambda} \mathbb{B} m_{\lambda}$ of symmetric functions with Boolean coefficients.

## $5 W^{\text {Sch }}(\mathbb{B})$

The purpose of this section is to determine $W^{\text {Sch }}(\mathbb{B})$. We will see in this section that it behaves very similarly to $W(\mathbb{B})$ in several ways. To stress these similarities, we are going to have the present section mirror Sect. 4 as closely as we can.

### 5.1 Schur constituents and Schur prime ideals

For any symmetric function $g \in \Lambda^{\text {Sch }}$, let us say the (Schur) constituents of $g$ are the partitions $v$ such that $g_{v} \neq 0$ in the unique decomposition $g=\sum_{\nu} g_{\nu} s_{v}, g_{v} \in \mathbb{N}$, where the $s_{\nu}$ are the Schur symmetric functions.

Let us say that a partition ideal $I$ is (Schur) prime if $I \neq P$ and whenever partitions $\lambda, \mu \in P$ have the property that every Schur constituent of $s_{\lambda} s_{\mu}$ is in $I$, then either $\lambda$ or $\mu$ lies in $I$. It will follow from (5.10) below that Schur primality is equivalent to monomial primality; the distinction between the two concepts is only needed until we prove that.

For any linear map $f \in \operatorname{Mod}_{\mathbb{N}}\left(\Lambda^{\text {Sch }}, \mathbb{B}\right)$, write

$$
\operatorname{pker}(f)=\left\{\lambda \in P \mid f\left(s_{\lambda}\right)=0\right\} .
$$

We will call pker $(f)$ the (Schur) partition kernel of $f$.

### 5.2 Constituents of products of Schur functions

The combinatorial lemmas about monomial constituents of $m_{\lambda} m_{\mu}$ that came to our aid in studying $W(\mathbb{B})$ have Schur analogues. These analogues rest upon the LittlewoodRichardson rule, which is a formula for the coefficient of $s_{\nu}$ in the product $s_{\lambda} s_{\mu}$. In order to be clear about the form of the rule we will use, we will fix some terms. For the basic terminology of tableaux, see Stanley's book [24, section 7.10, p. 309].

The reverse reading word of a tableau $T$ is defined to be the word obtained as follows: Write down the first row of $T$ in reverse order (that is, from right to left); after it, write down the second row of $T$ in reverse order; after it, the third one, and so on.

A word $w$ whose letters are positive integers is said to be a lattice permutation if and only if every positive integer $i$ and every $j \in\{0,1, \ldots, n\}$ (where $n$ is the length of $w$ ) satisfy the following condition: The number of $i$ 's among the first $j$ letters of $w$ is at least as large as the number of $i+1$ 's among the first $j$ letters of $w$.

A skew semistandard tableau $T$ is said to be a Littlewood-Richardson tableau if and only if its reverse reading word is a lattice permutation.

If a tableau $T$ is filled with positive integers, then the type of the tableau $T$ is defined to be the sequence ( $n_{1}, n_{2}, n_{3}, \ldots$ ), where $n_{k}$ is the number of boxes of $T$ filled with the integer $k$. Note that the type of a tableau $T$ is not necessarily a partition, but it is a partition whenever $T$ is a Littlewood-Richardson tableau.

For any partitions $\lambda, \mu$ and $\nu$, we let $c_{\lambda \mu}^{\nu}$ denote the number of LittlewoodRichardson tableaux of shape $\nu / \lambda$ and type $\mu$. When $\lambda \nsubseteq \nu$, this number is understood to be 0 . It is easy to see that $c_{\lambda \mu}^{\nu}=0$ unless $|\lambda|+|\mu|=|\nu|$. It is much less obvious that $c_{\lambda \mu}^{\nu}=c_{\mu \lambda}^{\nu}$ for all partitions $\lambda, \mu$ and $\nu$. The numbers $c_{\lambda \mu}^{\nu}$ are called the Littlewood-Richardson coefficients.

The Littlewood-Richardson rule then states

$$
\begin{equation*}
s_{\lambda} s_{\mu}=\sum_{\nu \in P} c_{\lambda \mu}^{\nu} s_{\nu} \tag{5.2.1}
\end{equation*}
$$

for any partitions $\lambda, \mu$. One can find a proof in Stanley's book [24, theorem A1.3.3, p. 432, appendix to Chapter 7]. For other points of view, see Sagan's book [20, theorem 4.9.4] or van Leeuwen [26, 1.4.4-1.4.5].

The following analogue of (4.2) is then clear.
Proposition 5.3 Let $\lambda$ and $\mu$ be two partitions. Then the Schur constituents of the product $s_{\lambda} s_{\mu}$ are precisely those partitions $v$ for which there exists a LittlewoodRichardson tableau of shape $\nu / \lambda$ and type $\mu$.

A corollary analogous to (4.3) results:
Corollary 5.4 Let $\lambda$ and $\mu$ be two partitions. Then, every Schur constituent vof $s_{\lambda} s_{\mu}$ satisfies $\lambda \subseteq \nu$ and $\mu \subseteq \nu$.

Proof Let $v$ be a Schur constituent of $s_{\lambda} s_{\mu}$. Due to (5.3), there exists a LittlewoodRichardson tableau of shape $\nu / \lambda$ and type $\mu$. As a consequence, the shape $\nu / \lambda$ must be well defined, so that $\lambda \subseteq \nu$. Similarly, $\mu \subseteq \nu$.

There is also a counterpart of (4.4) in the Schur function setting:
Proposition 5.5 Let $\lambda$ and $\mu$ be two partitions, and let $i$ and $j$ be two positive integers such that $\lambda_{i}>\mu_{i}, \lambda_{i}>\lambda_{i+1}, \mu_{j}>\lambda_{j}$, and $\mu_{j}>\mu_{j+1}$. Let $\bar{\lambda}$ denote the result of diminishing the ith part of $\lambda$ by 1 , and let $\bar{\mu}$ be the result of diminishing the $j$ th part of $\mu$ by 1. Then, for every Schur constituent $v$ of $s_{\bar{\lambda}} s_{\bar{\mu}}$, we have either $\lambda \subseteq v$ or $\mu \subseteq v$.

Proof Since $\lambda_{i}>\mu_{i}$ and $\mu_{j}>\lambda_{j}$, we have $i \neq j$. Thus, we can assume by symmetry that $i<j$.

For every partition $\gamma$, let $\lambda_{i}$ denote the result of removing the $i$ th part from the partition $\gamma$. By the definition of $\bar{\lambda}$, we have $\bar{\lambda} \backslash i=\lambda \backslash i$.

Fix a Schur constituent $\nu$ of $s_{\bar{\lambda}} s_{\bar{\mu}}$. Then, by (5.3), there exists a LittlewoodRichardson tableau of shape $\nu / \bar{\lambda}$ and type $\bar{\mu}$. Denote this tableau by $T$. Since $T$ exists, we have $\bar{\lambda} \subseteq v$. In particular, $\nu_{i} \geq \bar{\lambda}_{i} \geq \mu_{i} \geq \mu_{j}$ (since $i<j$ ).

First suppose that the $i$ th row of $\nu / \bar{\lambda}$ is nonempty. Then we have $\nu_{i} \geq \bar{\lambda}_{i}+1=\lambda_{i}$. Since $v$ is a Schur constituent of $s_{\bar{\lambda}} s_{\bar{\mu}}$, we have by (5.4) that $\bar{\lambda} \subseteq v$. Therefore we have $\lambda \subseteq \nu$, which finishes the proof in this case.

Now suppose that the $i$ th row of $\nu / \bar{\lambda}$ is empty. In this case, we will show $\mu \subseteq \nu$. Since we have $\bar{\mu} \subseteq \nu$, as above, it is enough to show $v_{j} \geq \mu_{j}$.

We can obtain a new tableau $T^{\prime}$ from the tableau $T$ by removing its (empty) $i$ th row and then moving each row below the $i$ th one up by one unit of length. This new tableau $T^{\prime}$ is a Littlewood-Richardson tableau of shape $\nu_{\backslash i} / \bar{\lambda} \backslash i$ and type $\bar{\mu}$ (since the original tableau $T$ was a Littlewood-Richardson tableau of shape $\nu / \bar{\lambda}$ and type $\bar{\mu}$ ). Hence, there exists a Littlewood-Richardson tableau of shape $\nu_{\backslash i} / \bar{\lambda} \backslash i$ and type $\bar{\mu}$. By (5.3) (applied to $\bar{\lambda} \backslash i, \bar{\mu}$ and $\nu_{\backslash i}$ instead of $\lambda, \mu$ and $\nu$ ), this yields that $\nu_{\backslash i}$ is a Schur constituent of the product $s_{\bar{\lambda}}^{\bar{\lambda}} s_{i} s_{\bar{\mu}}$. By (5.4) (again applied to $\bar{\lambda}_{\backslash i}, \bar{\mu}$ and $\nu_{\backslash i}$ instead of $\lambda, \mu$ and $\nu$ ), this yields that $\bar{\lambda} \backslash i \subseteq \nu_{\backslash i}$ and $\bar{\mu} \subseteq \nu_{\backslash i}$. But the partition $\bar{\mu}$ has at least $j-1$ parts at least $\mu_{j}$ (namely, its first $j-1$ parts, which as we know are unchanged from $\mu$ ). Since $\bar{\mu} \subseteq \nu_{\backslash i}$, this implies that the partition $\nu_{\backslash i}$ also has at least $j-1$ parts at least $\mu_{j}$. Hence, the partition $v$ has at least $j$ parts at least $\mu_{j}$ (because it contains all the parts of $\nu_{\backslash i}$, along with $\nu_{i}$ which as we know is at least $\mu_{j}$ ). In other words, $v_{j} \geq \mu_{j}$.

We will need two more facts about Schur constituents. The following is an analogue of (4.5). It will also be used only in Sect. 9.

Proposition 5.6 Let $\lambda$ and $\mu$ be two partitions.
(1) $\lambda+\mu$ is a Schur constituent of $s_{\lambda} s_{\mu}$.
(2) $\lambda \uplus \mu$ is a Schur constituent of $s_{\lambda} s_{\mu}$, where $\uplus$ is as in (4.5).
(3) If $\lambda$ and $\mu$ are both nonzero, then $s_{\lambda} s_{\mu}$ has at least two distinct nonzero Schur constituents.

Proof (1) By (5.3), it is enough to show that there exists a Littlewood-Richardson tableau of shape $(\lambda+\mu) / \lambda$ and type $\mu$. Such a tableau can be made by filling the $i$ th row of the diagram $(\lambda+\mu) / \lambda$ with the number $i$ for every positive integer $i$.
(2) For any partition $\gamma$, let $\gamma^{\prime}$ denote the conjugate partition. Since there is a $\mathbb{Z}$ algebra isomorphism $\omega: \Lambda_{\mathbb{Z}} \rightarrow \Lambda_{\mathbb{Z}}$ satisfying $\omega\left(s_{\gamma}\right)=s_{\gamma^{\prime}}$ for every partition $\gamma$, it
is enough to check that $(\lambda \uplus \mu)^{\prime}$ is a Schur constituent of $s_{\lambda^{\prime}} s_{\mu^{\prime}}$. But since we have $\lambda^{\prime}+\mu^{\prime}=(\lambda \uplus \mu)^{\prime}$, we are done by part (1).
(3) Since $\lambda$ and $\mu$ are nonzero, $\lambda+\mu$ and $\lambda \uplus \mu$ are distinct and nonzero. By parts (1) and (2), they are also Schur constituents of $s_{\lambda} s_{\mu}$.

### 5.7 Remark

When $\lambda$ and $\mu$ are two partitions, the two partitions $\lambda+\mu$ and $\lambda \uplus \mu$ are, respectively, the highest and the lowest Schur constituents of $s_{\lambda} s_{\mu}$ with respect to the so-called dominance order on partitions (also known as the majorization order)-the order in which a partition $\gamma$ is at least a partition $\delta$ of the same size if and only if every positive integer $i$ satisfies $\gamma_{1}+\gamma_{2}+\cdots+\gamma_{i} \geq \delta_{1}+\delta_{2}+\cdots+\delta_{i}$.

Proposition 5.8 Let $\lambda$ and $\mu$ be two partitions. Let $x$ be a positive integer. Let $\nu$ denote the partition obtained by sorting the components of the vector

$$
\left(\lambda_{1}+\mu_{1}, \ldots, \lambda_{x}+\mu_{x}, \lambda_{x+1}, \mu_{x+1}, \lambda_{x+2}, \mu_{x+2}, \ldots\right)
$$

in weakly decreasing order. Then $v$ is a Schur constituent of $s_{\lambda} s_{\mu}$.
Proof According to (5.3), it is enough to show that there exists a LittlewoodRichardson tableau of shape $\nu / \lambda$ and type $\mu$. We are going to construct such a tableau $T$.

Indeed, we are going to give two ways to construct such a $T$ —one explicit way which we will only sketch, and one less explicit one which we will detail.

First construction of $T$ : Let us first notice that $\nu_{i}=\lambda_{i}+\mu_{i}$ for every $i \in$ $\{1,2, \ldots, x\}$, so that the $i$ th row of $\nu / \lambda$ has length $\mu_{i}$ for every $i \in\{1,2, \ldots, x\}$. Let us fill this row with $i$ 's. Thus, the first $x$ rows of $\nu / \lambda$ are filled.

Now, for every cell $(i, j)$ of $v / \lambda$ with $i>x$, we fill the cell $(i, j)$ with $x$ plus the number of all $u \in\{x+1, x+2, \ldots, i\}$ satisfying $\lambda_{u}<j$.

This gives us a tableau $T$ of shape $\nu / \lambda$, which the reader can verify to be a Littlewood-Richardson tableau of type $\mu$ (the verification is not trivial, but does not require any new ideas). This proves the existence of such a tableau and thus concludes the proof of (5.8).

Second construction of $T$ : For every partition $\gamma$ and any integer $y$, let $\gamma \leq y$ denote the partition $\left(\gamma_{1}, \gamma_{2}, \ldots, \gamma_{y}\right)$, and let $\gamma_{>y}$ denote the partition $\left(\gamma_{y+1}, \gamma_{y+2}, \ldots\right)$. Clearly, $\nu_{\leq x}=\lambda_{\leq x}+\mu_{\leq x}$ and $\nu_{>x}=\lambda_{>x} \uplus \mu_{>x}$, where the operators + and $\uplus$ on partitions are defined as in (5.6).

Since $\nu_{\leq x}=\lambda_{\leq x}+\mu_{\leq x}$, the result (5.6)(1) (applied to $\lambda_{\leq x}$ and $\mu_{\leq x}$ instead of $\lambda$ and $\mu$ ) implies that $\nu_{\leq x}$ is a Schur constituent of $s_{\lambda \leq x} s_{\mu_{\leq x}}$. Hence, by (5.3), there exists a Littlewood-Richardson tableau of shape $\nu_{\leq x} / \lambda_{\leq x}$ and type $\mu_{\leq x}$. Let $T_{1}$ be such a tableau.

Since $\nu_{>x}=\lambda_{>x} \uplus \mu_{>x}$, the result (5.6)(2) (applied to $\lambda_{>x}$ and $\mu_{>x}$ instead of $\lambda$ and $\mu$ ) implies that $\nu_{>x}$ is a Schur constituent of $s_{\lambda_{>x}} s_{\mu_{>x}}$. Hence, by (5.3), there exists a Littlewood-Richardson tableau of shape $\nu_{>x} / \lambda_{>x}$ and type $\mu_{>x}$. Let
$T_{2}$ be such a tableau. Denote by $\widehat{T_{2}}$ the result of moving the tableau $T_{2}$ south (i.e., down) by $x$ rows and adding $x$ to each of its entries. It is easy to see that $\widehat{T_{2}}$ is a Littlewood-Richardson tableau of shape $v /\left(\nu_{1}, \nu_{2}, \ldots, \nu_{x}, \lambda_{x+1}, \lambda_{x+2}, \lambda_{x+3}, \ldots\right)$. Hence, $\widehat{T}_{2}$ doesn't intersect $T_{1}$, and moreover, each cell of $\widehat{T_{2}}$ lies strictly southwest of each cell of $T_{1}$. Thus, it is rather clear that overlaying $T_{1}$ with $\widehat{T_{2}}$ gives us a LittlewoodRichardson tableau of shape $\nu / \lambda$ and type $\mu$. This once again proves that such a tableau exists, and as we know this establishes (5.8).
Proposition 5.9 Let $2^{P}$ denote the set of subsets of $P$.
(1) The map $f \mapsto \operatorname{pker}(f)$ is a bijection $\operatorname{Mod}_{\mathbb{N}}\left(\Lambda^{\mathrm{Sch}}, \mathbb{B}\right) \rightarrow 2^{P}$.
(2) $\operatorname{pker}(f)$ is a partition ideal if and only if $f$ satisfies $f(x)=0 \Rightarrow f(x y)=0$.
(3) $\operatorname{pker}(f)$ is a Schur prime partition ideal if and only if $f$ is an $\mathbb{N}$-algebra map.

Proof Proceed as in (4.6), but invoke (5.3) and (5.4) instead of (4.2) and (4.3). (Notice that the $a_{\mu}$ in the proof of part (2) are now all equal to 1 according to the Pieri rule.)

It turns out that the Schur prime partition ideals are exactly the monomial prime partition ideals. Indeed, we have the following analogue of (4.7):

Proposition 5.10 For integers $x, y \in \mathbb{N}$, consider the partition ideal

$$
I_{(x, y)}=\left\{\lambda \in P \mid(y+1)^{x+1} \subseteq \lambda\right\}
$$

as defined in (4.7). Then $I_{(x, y)}$ is Schur prime. Conversely, every nonempty Schur prime partition ideal is of this form.

Proof This proof proceeds by adapting the proof of (4.7). The changes that are required (like replacing "prime" by "Schur prime," and referring to (5.5) and (5.4) instead of (4.4) and (4.3)) are almost all obvious; the only nontrivial change is to replace the reference to (4.2) by a reference to (5.8).

Proposition 5.11 Let $\infty$ denote the element of $W^{\text {Sch }}(\mathbb{B})$ satisfying $\infty\left(s_{\lambda}\right)=1$ for all partitions $\lambda$, and let $\{1\} \in W^{\text {Sch }}(\mathbb{B})$ be as in (2.3.3). Then we have
(1) $\operatorname{pker}(x+\{1\} y)=I_{(x, y)}$, for all $x, y \in \mathbb{N}$
(2) $\infty+z=\infty$, for all $z \in W^{S c h}(\mathbb{B})$
(3) $\infty \cdot z=\infty$, for all nonzero $z \in W^{\text {Sch }}(\mathbb{B})$
(4) $s_{\lambda}(\infty)=\infty$, for all nonzero partitions $\lambda$.

Proof (1) For clarity, given any $z \in \mathbb{N}$, let us write $f_{z}: \Lambda^{\text {Sch }} \rightarrow \mathbb{B}$ for the image of $z$ under the unique semiring map $\mathbb{N} \rightarrow W^{\operatorname{Sch}}(\mathbb{B})$. So if we let $I$ denote the partition ideal $\operatorname{pker}\left(f_{x}+\{1\} f_{y}\right)$, we need to show $I=I_{(x, y)}$.

In the case where $y=0$, the argument is identical to that for (4.9)(1). On the other hand, if $x=0$, then we have $\left(\{1\} f_{y}\right)\left(s_{\lambda}\right)=f_{y}\left(s_{\lambda^{\prime}}\right)$, where $\lambda^{\prime}$ is the conjugate partition of $\lambda$. (See for example [5, (6.11)].) And so this case follows from the previous one.

Now consider the general case. By (5.9) and (5.10), $I$ is either empty or of the form $I_{\left(x^{\prime}, y^{\prime}\right)}$. Therefore it is enough to show the three statements

$$
\begin{equation*}
(y+1)^{x+1} \in I, \quad y^{x+1} \notin I, \quad(y+1)^{x} \notin I . \tag{5.11.1}
\end{equation*}
$$

Since we have

$$
\Delta^{+}\left(s_{\lambda}\right)=\sum_{\mu, \nu} c_{\mu \nu}^{\lambda} s_{\mu} \otimes s_{\nu}
$$

by the definition of addition of Witt vectors, we have

$$
\left(f_{x}+\{1\} f_{y}\right)\left(s_{\lambda}\right)=\sum_{\mu, \nu} c_{\mu \nu}^{\lambda} \cdot f_{x}\left(s_{\mu}\right) \cdot\left(\{1\} f_{y}\right)\left(s_{\nu}\right)
$$

Therefore a partition $\lambda$ is in $I$ if and only if we have $\mu \in \operatorname{pker}\left(f_{x}\right)$ whenever $c_{\mu \nu}^{\lambda} \geq 1$ and $v \notin \operatorname{pker}\left(\{1\} f_{y}\right)$. Combining this with the two initial cases above we see that $\lambda$ is in $I$ if and only if the following implication holds:

$$
\left(c_{\mu \nu}^{\lambda} \geq 1 \text { and } \nu_{1} \leq y\right) \Longrightarrow \mu_{x+1} \geq 1
$$

We will now use this criterion to show each of the three statements in (5.11.1).
First consider $\lambda=(y+1)^{x+1}$, and assume $c_{\mu \nu}^{\lambda} \geq 1$ and $\nu_{1} \leq y$. Then there is a Littlewood-Richardson tableau of shape $\lambda / v$ and type $\mu$. Since $\nu_{1} \leq y$, the diagram of $\lambda / \nu$ contains the entire rightmost column of $\lambda$. Therefore all entries in that column of the tableau must be distinct, and hence $\mu$ has length at least $x+1$.

Now consider $\lambda=y^{x+1}$, and set $v=\lambda$ and $\mu=0$. Then there is clearly a Littlewood-Richardson tableau of shape $\lambda / v$ and type $\mu$, and further we have $\nu_{1} \leq y$ but $\mu_{x+1}=0$. Therefore $y^{x+1} \notin I$. Similarly, for $\lambda=(y+1)^{x}$, take $\nu=0$ and $\mu=\lambda$. We have $\nu_{1} \leq y$, and there is clearly a Littlewood-Richardson tableau of shape $\lambda / \nu=\lambda$ and type $\mu=\lambda$, but $\mu_{x+1}=0$. Therefore $(y+1)^{x} \notin I$.
(2) Since $\infty+z \geq \infty$, we have $\infty+z \succcurlyeq \infty$, by (3.4), and hence $\infty+z=\infty$.
(3) By distributivity and parts (1) and (2), it is enough to check this for the $\mathbb{N}$ module generators $z=1,\{1\}, \infty$. It is clear for $z=1$. For $z=\{1\}$, it holds because multiplication by $\{1\}$ is an involution which preserves $W^{\text {Sch }}(\mathbb{B}) \backslash\{\infty\}$, by part (1), and hence fixes $\infty$. When $z=\infty$, the argument is that of (4.9)(2) but with the Schur basis instead of the monomial one.
(4) By definition $\left(s_{\lambda}(\infty)\right)\left(s_{\mu}\right)=\infty\left(s_{\mu} \circ s_{\lambda}\right)$. So the statement to be shown is equivalent to the statement $s_{\mu} \circ s_{\lambda} \neq 0$ for all partitions $\mu$. Since $\lambda \neq 0$, we can write $s_{\lambda}=y_{1}+y_{2}+\cdots$, where each $y_{i}$ is a monomial in the formal variables $x_{1}, x_{2}, \ldots$ with coefficient 1 . Therefore we have $s_{\mu} \circ s_{\lambda}=s_{\mu}\left(y_{1}, y_{2}, \ldots\right)$, but this is nonzero since all Schur functions are nonzero.

### 5.12 Explicit description of $W^{\text {Sch }}(\mathbb{B})$

Here we prove part (3) of Theorem (1.1), although we use slightly different notation. It follows from (5.9), (5.10) and (5.11)(1) that we have a bijection

$$
\begin{equation*}
\mathbb{N}[\eta] /\left(\eta^{2}=1\right) \cup\{\infty\} \xrightarrow{\sim} W^{\text {Sch }}(\mathbb{B}) \tag{5.12.1}
\end{equation*}
$$

sending $x+y \eta$ to $x+\{1\} y$ and sending $\infty$ to $\infty$. It is an isomorphism of $\mathbb{N}$-algebras if we give the left-hand side the $\mathbb{N}$-algebra structure extending that on $\mathbb{N}[\eta] /\left(\eta^{2}=1\right)$ and satisfying $\infty+z=\infty$ and $\infty \cdot z=\infty$, for $z \neq 0$. This follows from (5.11) and the equality $\{1\}^{2}=1$. It is also clearly order preserving, as in (4.9)(7). The $\Lambda^{\text {Sch }}$-semiring structure on $\mathbb{N}[\eta] /\left(\eta^{2}=1\right) \cup\{\infty\}$ induced by this isomorphism is the unique one satisfying $s_{\lambda}(\infty)=\infty$ and

$$
s_{\lambda}(\eta)= \begin{cases}\eta^{r} & \text { if } \lambda=\left(1^{r}\right) \\ 0 & \text { otherwise }\end{cases}
$$

Proposition 5.13 In terms of the bijections (4.8.1) and (5.12.1), the map $W(\mathbb{B}) \rightarrow$ $W^{\text {Sch }}(\mathbb{B})$ sends any element of the form $(x, 0)$ to $x=x+\{1\} 0$. It sends all other elements to $\infty$.

Proof $\operatorname{By}(4.9)(1)$, the pair $(x, 0)$ is the image of $x$ under the unique map $\mathbb{N} \rightarrow W(\mathbb{B})$; so it must map to $x$ in $W^{S c h}(\mathbb{B})$. Therefore all that remains is to show that the preimage of $\infty$ under the map $W(\mathbb{B}) \rightarrow W^{\text {Sch }}(\mathbb{B})$ contains all other elements. By algebraic laws listed in Theorem (1.1), it is enough to show the preimage of $\infty$ contains $\infty$ and all elements of the form $(0, y)$, where $y \geq 1$.

Let $\lambda$ be a partition, and put $r=\sum_{i} \lambda_{i}$. Then it is easily seen that there is at least one semistandard Young tableau of shape $\lambda$ and type $1^{r}$. So the Kostka number $K_{\lambda, 1^{r}}$ is at least 1 . Therefore we have $s_{\lambda}=m_{1^{r}}+\cdots \geq m_{1^{r}}$, and hence $(0, y)\left(s_{\lambda}\right) \geq$ $(0, y)\left(m_{1}{ }^{r}\right)=1$ for any $y \geq 1$. Therefore any such $(0, y)$ maps to $\infty$ in $W^{\operatorname{Sch}}(\mathbb{B})$. Last, $\infty$ also maps to $\infty$, because the map $\infty: \Lambda_{\mathbb{N}} \rightarrow \mathbb{B}$ is 1 on all nonzero elements of $\Lambda_{\mathbb{N}}$ and therefore on all Schur functions $s_{\lambda}$.

### 5.14 Remark

Observe that while the map $W(A) \rightarrow W^{\text {Sch }}(A)$ is injective when $A$ has additive cancellation, it is not so in general. Indeed, this fails for $A=\mathbb{B}$.

## 6 Total positivity

In this section, we recall the relationship between the Schur Witt vectors and Schoenberg's theory of total positivity. It gives another technique for understanding Witt vectors of $\mathbb{R}_{+}$and hence of semirings that map to $\mathbb{R}_{+}$. All details can be found in section 7 of [5].

### 6.1 Total positivity

As in the introduction, consider the map

$$
\begin{align*}
\sigma: W^{\mathrm{Sch}}(A) & \longrightarrow 1+t A[[t]]  \tag{6.1.1}\\
a & \longmapsto 1+a\left(e_{1}\right) t+a\left(e_{2}\right) t^{2}+\cdots,
\end{align*}
$$

where the $e_{n}$ are the elementary symmetric functions

$$
e_{n}=\sum_{i_{1}<\cdots<i_{n}} x_{i_{1}} \cdots x_{i_{n}} \in \Lambda^{\mathrm{Sch}}
$$

It is a bijection if $A$ is a ring and an injection if $A$ has additive cancellation.
Following Schoenberg [22], one says that a formal real series $f(t)=\sum_{i \in \mathbb{Z}} a_{i} t^{i}$ is totally positive if all the minors of the infinite matrix $\left(a_{i-j}\right)_{i, j}$ are $\geq 0$. The connection between total positivity and $W^{\text {Sch }}$ is given by two standard facts in the theory of symmetric functions, that $\Lambda^{\text {Sch }}$ agrees with the $\mathbb{N}$-linear span of the skew Schur functions and that the skew Schur functions are the minors of the matrix $\left(e_{i-j}\right)_{i, j}$. It follows that for any subring $B \subseteq \mathbb{R}$, the map (6.1.1) restricts to a bijection

$$
\sigma: W^{\mathrm{Sch}}\left(B \cap \mathbb{R}_{+}\right) \xrightarrow{\sim}\{\text { totally positive series in } 1+t B[[t]]\} .
$$

The classification of the totally positive series given by the Edrei-Thoma theorem then amounts to the following:

Theorem 6.2 The map $\sigma: W(\mathbb{R}) \rightarrow 1+t \mathbb{R}[[t]]$ restricts to a bijection

$$
\begin{equation*}
W^{\mathrm{Sch}}\left(\mathbb{R}_{+}\right) \xrightarrow{\sim}\left\{\left.e^{\gamma t} \frac{\prod_{i=1}^{\infty}\left(1+\alpha_{i} t\right)}{\prod_{i=1}^{\infty}\left(1-\beta_{i} t\right)} \right\rvert\, \alpha_{i}, \beta_{i}, \gamma \in \mathbb{R}_{+}, \sum_{i}\left(\alpha_{i}+\beta_{i}\right)<\infty\right\} \tag{6.2.1}
\end{equation*}
$$

This was conjectured by Schoenberg [22] and proved by Aissen-SchoenbergWhitney [2] and, in the final step, Edrei [1,10]. It was independently discovered and proved by Thoma in his work [25] on the asymptotic character theory of symmetric groups. Both arguments use hard results in the theory of entire functions. Later, proofs entirely within asymptotic character theory emerged out of the work of Vershik [27], Kerov [15], Okounkov [19] and Olshanski [16].

The following is an easy consequence of (6.2) (see section 7 of [5]):
Corollary 6.3 The map $\sigma: W(\mathbb{R}) \rightarrow 1+t \mathbb{R}[[t]]$ restricts to bijections

$$
\begin{aligned}
& W\left(\mathbb{R}_{+}\right) \xrightarrow{\sim}\left\{e^{\gamma t} \prod_{i=1}^{\infty}\left(1+\alpha_{i} t\right) \mid \alpha_{i}, \gamma \in \mathbb{R}_{+}, \sum_{i} \alpha_{i}<\infty\right\} \\
& W(\mathbb{N}) \xrightarrow{\sim}\{f(t) \in 1+t \mathbb{Z}[t] \mid \text { all complex roots of } f(t) \text { are negative reals }\}
\end{aligned}
$$

The algebra $W^{S c h}(\mathbb{N})$ is determined in the following section.

## 7 Edrei-Thoma and Boolean Witt vectors

The purpose of this section is to describe the map on Witt vectors induced by the map $\mathbb{R}_{+} \rightarrow \mathbb{B}$ in terms of the power-series description of $W\left(\mathbb{R}_{+}\right)$and the combinatorial description of $W(\mathbb{B})$. We then use this to prove Theorem (1.2), or equivalently (1.3).

Proposition 7.1 (1) In terms of (6.3) and (1.1), the functorial map $W\left(\mathbb{R}_{+}\right) \rightarrow W(\mathbb{B})$ satisfies

$$
e^{\gamma t} f(t) \longmapsto \begin{cases}(\operatorname{deg} f(t), 0) & \text { if } \gamma=0  \tag{7.1.1}\\ (\operatorname{deg} f(t), 1) & \text { if } \gamma \neq 0\end{cases}
$$

for any polynomial $f(t)$. It sends all other series to $\infty$.
(2) In terms of (6.2) and (1.1), the functorial map $W^{\text {Sch }}\left(\mathbb{R}_{+}\right) \rightarrow W^{\text {Sch }}(\mathbb{B})$ sends any rational function $f(t) / g(t)$ in lowest terms to the pair $(\operatorname{deg} f(t), \operatorname{deg} g(t))$. It sends all other series to $\infty$.

Proof (1) Let $\varphi$ denote the functorial map $\varphi: W\left(\mathbb{R}_{+}\right) \rightarrow W(\mathbb{B})$. Let $z \in W\left(\mathbb{R}_{+}\right)$be a Witt vector. By (6.3), we can write

$$
\sigma(z)=e^{\gamma t} \prod_{i=1}^{\infty}\left(1+\alpha_{i} t\right)
$$

Let $r \leq \infty$ be the number of nonzero $\alpha_{i}$. Abusing notation, let us still write $r$ for the image of $r$ under the map $\mathbb{N} \cup\{\infty\} \rightarrow W(\mathbb{B})$. (So $r=(r, 0)$ for $r \neq \infty$.) By (4.9), it is enough to show $\varphi(z)=r$ if $\gamma=0$, and $\varphi(z)=r+(0,1)$ if $\gamma \neq 0$. It is also enough to consider separately the cases where $\gamma=0$ and where all $\alpha_{i}=0$; this is because the map $\sigma$ identifies addition of Witt vectors with multiplication of power series.

First suppose $\gamma=0$. Then we have $z\left(m_{\lambda}\right)=m_{\lambda}\left(\alpha_{1}, \alpha_{2}, \ldots\right)$ for all partitions $\lambda$. (For example, one can observe that $z$ is the image of $\sum_{i}\left[\alpha_{i}\right]$ under the map $\mathbb{N}\left[\mathbb{R}_{+}\right]^{\wedge} \rightarrow$ $W\left(\mathbb{R}_{+}\right)$from the convergent monoid algebra. See (7.6) and (4.12.1) in [5].) If $r=\infty$, this is zero for no $\lambda$, and if $r<\infty$, it is zero if and only if $\lambda_{r+1} \geq 1$. Therefore the partition ideal corresponding to $\varphi(z)$ is the empty one if $r=\infty$ and is $I_{(r, 0)}$ if $r<\infty$. Hence we have $\varphi(z)=r$, whether $r$ is infinite or not.

Now suppose $\gamma \neq 0$ and $\alpha_{i}=0$ for all $i$. Then we have

$$
\exp (\gamma t)=\sigma(z)=\exp \left(-\sum_{n \geq 1} z\left(\psi_{n}\right) \frac{(-t)^{n}}{n}\right)
$$

and hence $z\left(m_{1^{1}}\right)=z\left(\psi_{1}\right)=\gamma \neq 0$ and $z\left(m_{2^{1}}\right)=z\left(\psi_{2}\right)=0$. Therefore $\varphi(z)$ has partition kernel $I_{(0,1)}$, and so $\varphi(z)=(0,1)$.
(2) Let $\varphi^{\text {Sch }}$ denote the functorial map $W^{\text {Sch }}\left(\mathbb{R}_{+}\right) \rightarrow W^{\text {Sch }}(\mathbb{B})$. By (6.2) and (6.3), any element $z \in W^{\text {Sch }}\left(\mathbb{R}_{+}\right)$is of the form $x+\{1\} y$ for some $x, y \in W\left(\mathbb{R}_{+}\right)$. Because anti-Teichmüller lifts are functorial, we have $\varphi^{\operatorname{Sch}}(x+\{1\} y)=\varphi^{\text {Sch }}(x)+\{1\} \varphi^{\operatorname{Sch}}(y)$. Therefore by (5.11), we are reduced to the case $z \in W\left(\mathbb{R}_{+}\right)$, and this follows from (5.13) and part (1) above.

### 7.2 Remark

The Edrei-Thoma Theorem (6.2) provides a discrete invariant associated to a totally positive series, namely the pair consisting of the number of zeros and the number
of poles, or $\infty$ if the series is not a rational function. By (7.1)(2), we may interpret this invariant as the image of the corresponding Witt vector under the natural map $W^{\text {Sch }}\left(\mathbb{R}_{+}\right) \rightarrow W^{\text {Sch }}(\mathbb{B})$. So we have the satisfying fact that the coarse invariant coming from the Boolean Witt vectors discussed in the introduction agrees with the obvious discrete invariant given by the Edrei-Thoma theorem.

### 7.3 Remark

The map $W^{\text {Sch }}\left(\mathbb{R}_{+}\right) \rightarrow W^{\text {Sch }}(\mathbb{B})$ is surjective, but $W\left(\mathbb{R}_{+}\right) \rightarrow W(\mathbb{B})$ is not.
Lemma 7.4 Given any $a \in W^{\text {Sch }}\left(\mathbb{R}_{+}\right)$, we have $\lim _{k \rightarrow \infty} \min _{\lambda \vdash k} a\left(s_{\lambda}\right)=0$, where $\lambda \vdash k$ means that $\lambda$ is a partition of $k$.

Proof Let $k \in \mathbb{N}$. By Schur-Weyl duality, every vector space $V$ over $\mathbb{C}$ satisfies

$$
V^{\otimes k} \cong \bigoplus_{\lambda \vdash k} M_{\lambda} \otimes L_{\lambda}(V)
$$

as $S_{k} \times \mathrm{GL}(V)$-modules, where $M_{\lambda}$ is the irreducible $S_{k}$-module corresponding to the partition $\lambda$, and $L_{\lambda}$ is the Schur functor corresponding to the partition $\lambda$. Taking the GL( $V$ )-character of this isomorphism, we obtain

$$
s_{1}^{k}=\sum_{\lambda} \operatorname{dim}\left(M_{\lambda}\right) \cdot s_{\lambda},
$$

where we sum over partitions $\lambda$ satisfying $\sum_{i} \lambda_{i}=k$. Applying the ring homomorphism $a$ to this equality, we obtain

$$
\begin{equation*}
a\left(s_{1}\right)^{k}=\sum_{\lambda} \operatorname{dim}\left(M_{\lambda}\right) \cdot a\left(s_{\lambda}\right) \tag{7.4.1}
\end{equation*}
$$

Writing $c_{k}=\min _{\lambda \vdash k} a\left(s_{\lambda}\right)$, we then have

$$
a\left(s_{1}\right)^{k} \geq c_{k} \sum_{\lambda} \operatorname{dim}\left(M_{\lambda}\right) \geq c_{k}\left(\sum_{\lambda} \operatorname{dim}\left(M_{\lambda}\right)^{2}\right)^{1 / 2}=c_{k} \sqrt{k!},
$$

since the sum of the squares of the dimensions of all irreducible $S_{k}$-modules is $\left|S_{k}\right|=$ $k!$. It then follows that $c_{k} \rightarrow 0$ as $k \rightarrow \infty$.

### 7.5 Proof of Theorem (1.2)

For any Witt vector $a \in W^{\text {Sch }}(\mathbb{N})$, Lemma (7.4) implies there must be a partition $\lambda$ such that $a\left(s_{\lambda}\right)=0$. In other words, the image of $a$ in $W^{\text {Sch }}(\mathbb{B})$ is not $\infty$. Therefore by $(7.1)(2)$ and the Edrei-Thoma Theorem (6.2), the series $\sigma(a) \in 1+t \mathbb{R}[[t]]$ is of the form $g(t) / h(t)$, where $g(t)$ and $h(t)$ are polynomials in $1+t \mathbb{R}[t]$ such that all the
complex roots of $g(t)$ are negative real numbers and all those of $h(t)$ are positive real numbers. Now invoke the fact that if the ratio of two coprime polynomials in $1+t \mathbb{C}[t]$ has integral coefficients when expanded as a power series in $t$, then the polynomials themselves have integral coefficients. (To show they have rational coefficients, one can observe that $g(t)$ and $h(t)$ are invariant under any field automorphism of $\mathbb{C}$, but there are also choice-free arguments. Then to show the coefficients are integral, one can argue as in the solution to exercise 1a of chapter 4 of [23].)

### 7.6 Remark

Another way of expressing Theorem (1.2) is that we have an isomorphism of $\mathbb{N}$ algebras

$$
W(\mathbb{N}) \otimes_{\mathbb{N}} \mathbb{N}[\eta] /\left(\eta^{2}=1\right) \xrightarrow{\sim} W^{\operatorname{Sch}}(\mathbb{N})
$$

which is the natural map on the factor $W(\mathbb{N})$ and sends $1 \otimes \eta$ to the anti-Teichmüller $\operatorname{lift}\{1\} \in W^{\text {Sch }}(\mathbb{N})$. In fact, each factor on the left side has a compatible $\Lambda^{\text {Sch }}$-semiring structure making the map an isomorphism of $\Lambda^{\text {Sch }}$-semirings.

### 7.7 Remark: partition-ideal phenomena in general $W(A)$

As discussed in the introduction, phenomena related to $\mathbb{N}^{2}$ and partition ideals will be present in $W(A)$ (and $W^{\text {Sch }}(A)$ ) whenever $A$ is not a ring. This is simply because, by (7.9) below, every nonring $A$ admits a map to $\mathbb{B}$ and hence there is always a map $W(A) \rightarrow W(\mathbb{B})$. This is analogous to the presence of $p$-adic phenomena in Witt vectors of rings admitting a map to a field of characteristic $p$, or equivalently rings not containing $1 / p$. But in the case of $\mathbb{B}$, there is more structure: The set of maps $A \rightarrow \mathbb{B}$ has a partial order given by pointwise application of the partial order $0<1$ of $\mathbb{B}$. So in fact there is a family of partition-ideal phenomena indexed by this partially ordered set.

For example, consider the case where $A$ is a zerosumfree domain (that is, it is nonzero and satisfies the implication $x+y=0 \Rightarrow x=y=0$ and the implication $x y=0 \Rightarrow x=0$ or $y=0$ ) and $f$ is the $\mathbb{N}$-algebra map $A \rightarrow \mathbb{B}$ sending all nonzero elements to 1 . This is the unique maximal element of the partially ordered set and hence gives a particularly natural way in which partition-ideal phenomena appear. The following is an immediate a consequence of (1.1).

Corollary 7.8 Let A be a zerosumfree domain, and let $a \in W(A)$ be a Witt vector such that $a\left(m_{\lambda}\right)=0$ for some $\lambda \in P$. Then there is a unique pair $(x, y) \in \mathbb{N}^{2}$ such that $a\left(m_{\lambda}\right)=0$ if and only if $\lambda_{x+1} \geq y+1$. The analogous statement is true for $W^{\text {Sch }}(A)$ and the Schur basis $\left(s_{\lambda}\right)_{\lambda \in P}$.

We conclude this section with the proposition invoked above.
Proposition 7.9 (1) An $\mathbb{N}$-module $M$ is a group if and only if $\mathbb{B} \otimes_{\mathbb{N}} M=0$.
(2) An $\mathbb{N}$-algebra $A$ is a ring if (under the axiom of choice) and only if there exists no $\mathbb{N}$-algebra homomorphism $A \rightarrow \mathbb{B}$.

Proof (1) If $M$ is a group, then it is clear that $\mathbb{B} \otimes_{\mathbb{N}} M=0$. Conversely, suppose $\mathbb{B} \otimes_{\mathbb{N}} M=0$. Since $\mathbb{B}$ is $\mathbb{N} /(2=1)$, the module $\mathbb{B} \otimes_{\mathbb{N}} M$ is the quotient of $M$ by the equivalence relation $\approx$ generated by relations $2 x=x$ for all $x \in M$. That is, $\approx$ is the transitive closure of the reflexive symmetric relation $\sim$ defined by

$$
a \sim b \Longleftrightarrow a=z+2 x+y, b=z+x+2 y \text { for some } x, y, z \in M .
$$

Now the set of additively invertible elements of $M$ is closed under taking both sums and summands; so it follows that if $a \sim b$, then $a$ is additively invertible if and only if $b$ is additively invertible. Since $\approx$ is the transitive closure of $\sim$, the same is true of $\approx$. Since we have assumed $M / \approx$ is 0 , we have $a \approx 0$ for any element $a \in M$. Thus $a$ is additively invertible, and so $M$ is a group.
(2) If $A$ is a ring, then it is clear there are no maps $A \rightarrow \mathbb{B}$. Conversely, if $A$ is not a ring, then by part (1), we have $\mathbb{B} \otimes_{\mathbb{N}} A \neq 0$, and so by Zorn's lemma $\mathbb{B} \otimes_{\mathbb{N}} A$ has a minimal nonzero quotient $\mathbb{N}$-algebra $B$. Since $B$ is nonzero and has no nonzero quotients, it must be a field or $\mathbb{B}$. (See Golan [11, p. 87].) But since $B$ admits a map from $\mathbb{B}$, it cannot be a field. Therefore we have $B=\mathbb{B}$, and hence there is a map $A \rightarrow \mathbb{B}$.

## $8 W_{(p)}(\mathbb{B})$

Let $p$ be a prime number. In this section, we determine the $p$-typical Witt vectors of $\mathbb{B}$. The result is straightforward, but we include it for completeness.

## $8.1 p$-typical Witt vectors

Let us recall some definitions and basic results from section 8 of [5]. Write

$$
\psi=\sum_{i} x_{i}^{p}, \quad d=\frac{\left(\sum_{i} x_{i}\right)^{p}-\sum_{i} x_{i}^{p}}{p}, \quad d_{i, j}=d^{\circ i} \circ \psi^{\circ j} \in \Lambda_{\mathbb{N}},
$$

for all $i, j \in \mathbb{N}$, where $\circ$ denotes the plethysm operation. The $\mathbb{N}$-algebra $\Lambda_{\mathbb{N},(p)}$ of $p$-typical symmetric functions is then defined to be the sub- $\mathbb{N}$-algebra of $\Lambda_{\mathbb{N}}$ generated by the set $\left\{d_{i, j} \mid i, j \in \mathbb{N}\right\}$, and the set of $p$-typical Witt vectors with entries in a given $\mathbb{N}$-algebra $A$ is defined by

$$
W_{(p)}(A)=\operatorname{Alg}_{\mathbb{N}}\left(\Lambda_{\mathbb{N},(p)}, A\right)
$$

By (8.3) of [5], the relations $d_{i, j}^{p}=p d_{i+1, j}+d_{i, j+1}$ generate all the relations between the $d_{i, j}$. Therefore if we write $a_{i, j}=a\left(d_{i, j}\right)$, we have the identification

$$
\begin{equation*}
W_{(p)}(A)=\left\{\left(a_{i, j}\right) \in A^{\mathbb{N}^{2}} \mid a_{i, j}^{p}=p a_{i+1, j}+a_{i, j+1}\right\} . \tag{8.1.1}
\end{equation*}
$$

As shown in [5], there is a unique structure of a composition $\mathbb{N}$-algebra on $\Lambda_{\mathbb{N},(p)}$ compatible with that of $\Lambda_{\mathbb{N}}$. As always, $W_{(p)}(A)$ is then naturally an $\mathbb{N}$-algebra with an action of $\Lambda_{\mathbb{N},(p)}$ defined by $f(a): g \mapsto a(g \circ f)$ for all $f, g \in \Lambda_{\mathbb{N},(p)}$ and $a \in W_{(p)}(A)$. Therefore, in terms of (8.1.1), we have

$$
\begin{equation*}
d(a)_{i, j}=a_{i+1, j}, \quad \psi(a)_{i, j}=a_{i, j+1} \tag{8.1.2}
\end{equation*}
$$

## 8.2 $W_{(p)}(\mathbb{B})$

In the particular case $A=\mathbb{B}$, Eq. (8.1.1) simplifies to

$$
\begin{equation*}
W_{(p)}(\mathbb{B})=\left\{\left(a_{i, j}\right) \in \mathbb{B}^{\mathbb{N}^{2}} \mid a_{i, j}=0 \Leftrightarrow a_{i, j+1}=a_{i+1, j}=0\right\} . \tag{8.2.1}
\end{equation*}
$$

Since any $a_{i, j}$ is either 0 or 1 , we see that $W_{(p)}(\mathbb{B})$ is identified with the set of subsets $I \subseteq \mathbb{N}^{2}$ satisfying $(i, j) \in I \Longleftrightarrow(i+1, j),(i, j+1) \in I$. Since any such subset $I$, if nonempty, is determined by its unique element $(i, j)$ with $i+j$ minimal, we have a bijection

$$
\begin{equation*}
\mathbb{N}^{2} \cup\{\infty\} \xrightarrow{\sim} W_{(p)}(\mathbb{B}) \tag{8.2.2}
\end{equation*}
$$

sending $(x, y)$ to the Witt vector $\left(a_{i, j}\right)$ with $a_{i, j}=0 \Leftrightarrow i \geq x, j \geq y$ and sending $\infty$ to the Witt vector with $a_{i, j}=1$ for all $i, j$. By (3.2.3), this bijection is an isomorphism of partially ordered sets, where $W_{(p)}(\mathbb{B})$ has the partial order $\preccurlyeq$, where $\mathbb{N}^{2}$ has the componentwise order

$$
(x, y) \leq\left(x^{\prime}, y^{\prime}\right) \Longleftrightarrow x \leq x^{\prime} \text { and } y \leq y^{\prime}
$$

and where $\infty$ is a terminal element. Under this bijection, (8.1.2) becomes

$$
\begin{equation*}
d(x, y)=(x-1, y), \quad \psi(x, y)=(x, y-1), \quad d(\infty)=\infty, \quad \psi(\infty)=\infty \tag{8.2.3}
\end{equation*}
$$

where we understand that negative coordinates are rounded up to zero.
Proposition 8.3 The $\mathbb{N}$-algebra structure on $\mathbb{N}^{2} \cup\{\infty\}$ inherited from (8.2.2) satisfies the following properties for all $x, x^{\prime}, y, y^{\prime} \in \mathbb{N}$ :
(1) The image of $n \in \mathbb{N}$ under the unique algebra map $\mathbb{N} \rightarrow \mathbb{N}^{2} \cup\{\infty\}$ is $(n, 0)$ if $n \leq 1$ (or rather $n \leq 2$ if $p=2$ ) and is $\infty$ otherwise.
(2) $(x, 0)+\left(x^{\prime}, 0\right)=\infty$ if $x, x^{\prime} \geq 2$
(3) $(0, y)+\left(0, y^{\prime}\right)=\left(0, \max \left\{y, y^{\prime}\right\}\right)$
(4) $(x, 0)+(0, y)=(x, y)$
(5) $\infty+z=\infty$ for all $z \in \mathbb{N}^{2} \cup\{\infty\}$
(6) $(x, 0)\left(x^{\prime}, 0\right)=\infty$ for $x, x^{\prime} \geq 2$
(7) $(0, y)\left(0, y^{\prime}\right)=\left(0, \min \left\{y, y^{\prime}\right\}\right)$
(8) $(x, 0)(0, y)=(0, y)$ if $x \geq 1$
(9) $\infty(0, y)=(0, y)$
(10) $\infty(x, 0)=\infty$ if $x \geq 1$.

Proof We will identify any element of $\mathbb{N}^{2} \cup\{\infty\}$ with its image in $W_{(p)}(\mathbb{B})$ under the bijection (8.2.2).
(1) Let $\varphi$ denote the map $\mathbb{N} \rightarrow W_{(p)}(\mathbb{N})$. Since it is $\Lambda_{\mathbb{N},(p)}$-equivariant, we have

$$
(\varphi(n))_{i, 0}=\left(d^{\circ i}(\varphi(n))\right)_{0,0}=\left(\varphi\left(d^{\circ i}(n)\right)\right)_{0,0}=d^{\circ i}(n) .
$$

Now for $n \in \mathbb{N}$, we have $d(n)=\left(n^{p}-n\right) / p$. Thus for $p \geq 3$, we have $d(0)=d(1)=0$ and $d(m) \geq 2$ for $m \geq 2$. Therefore if $n \leq 1$, the smallest $i$ such that $d^{\circ i}(n)=0$ is $n$, and no such $i$ exists if $n \geq 2$. This finishes the proof for $p \geq 3$. The statement for $p=2$ follows similarly from $d(0)=d(1)=0, d(2)=1$, and $d(m) \geq 3$ for $m \geq 3$.
(2) By part (1), we have $(x, 0)+\left(x^{\prime}, 0\right) \succcurlyeq(2,0)+(2,0)=2+2=\infty$.
(3) Since the power sum $d_{0, j}=\psi^{\circ j}$ is primitive for the co-addition map $\Delta^{+}$, we have $\left((0, y)+\left(0, y^{\prime}\right)\right)_{0, j}=(0, y)_{0, j}+\left(0, y^{\prime}\right)_{0, j}$. This evaluates to 0 if and only if $y, y^{\prime} \leq j$. Therefore $\left((0, y)+\left(0, y^{\prime}\right)\right)_{i, j}$ is zero for $(i, j)=\left(0, \max \left\{y, y^{\prime}\right\}\right)$ but not for any smaller $(i, j)$. Thus we have $(0, y)+\left(0, y^{\prime}\right)=\left(0, \max \left\{y, y^{\prime}\right\}\right)$.
(4) If either $x$ or $y$ is zero, the statement follows because $(0,0)$ is the additive identity, by part (1). So assume $x, y \geq 1$. By (8.2.3), it is enough to show $\psi^{\circ y-1}((x, 0)+$ $(0, y))=(x, 1)$. Because we have

$$
\psi^{\circ y-1}((x, 0)+(0, y))=(x, 0)+(0,1)
$$

it is enough to show $(x, 0)+(0,1)=(x, 1)$ for $x \geq 1$. Since we have

$$
\psi((x, 0)+(0,1))=\psi(x, 0)+\psi(0,1)=(x, 0)+(0,0)=(x, 0)
$$

$(x, 0)+(0,1)$ is either $(x, 0)$ or $(x, 1)$, again by (8.2.3). To rule out $(x, 0)$, we can show

$$
d((x, 0)+(0,1)) \succcurlyeq(x-1,1) .
$$

This holds because we have

$$
\Delta^{+}(d)=d \otimes 1+1 \otimes d+\cdots
$$

and hence

$$
\begin{aligned}
d((x, 0)+(0,1)) & =d(x, 0)+d(0,1)+\cdots \\
& \succcurlyeq d(x, 0)+d(0,1) \\
& =(x-1,0)+(0,1) \\
& =(x-1,1)
\end{aligned}
$$

by (3.4) and induction on $x$.
(5) By (3.4), we have $\infty+z \succcurlyeq \infty$ and hence $\infty+z=\infty$.
(6) Since $\succcurlyeq$ is an Alg $_{\mathbb{N}}$-preorder, we have $(x, 0)\left(x^{\prime}, 0\right) \succcurlyeq(2,0)(2,0)$. Therefore it is enough to show $(2,0)^{2}=\infty$. We have

$$
\begin{equation*}
\Delta^{\times}(d)=\psi \otimes d+d \otimes \psi+p d \otimes d \tag{8.3.1}
\end{equation*}
$$

and hence

$$
\begin{aligned}
d((2,0)(2,0)) & =\psi(2,0) d(2,0)+d(2,0) \psi(2,0)+p d(2,0)^{2} \\
& =(2,0)(1,0)+(1,0)(2,0)+p(1,0)^{2} \\
& =(2,0)+(2,0)+p(1,0) \\
& =\infty+p(1,0) \\
& =\infty
\end{aligned}
$$

by (8.2.3) and parts (1), (2) and (5). It follows that $(0,2)^{2}=\infty$.
(7) For any $j$, the element $d_{0, j}=\psi^{\circ j}$ is group-like under $\Delta^{\times}$. Therefore we have

$$
\left((0, y)\left(0, y^{\prime}\right)\right)_{0, j}=((0, y))_{0, j} \cdot\left(\left(0, y^{\prime}\right)\right)_{0, j}
$$

which is 0 if and only if $j \geq \min \left\{y, y^{\prime}\right\}$. Therefore $\left((0, y)\left(0, y^{\prime}\right)\right)_{i, j}$ is zero for $(i, j)=\left(0, \min \left\{y, y^{\prime}\right\}\right)$ but not for any smaller $(i, j)$. It follows that $(0, y)\left(0, y^{\prime}\right)$ is $\left(0, \min \left\{y, y^{\prime}\right\}\right)$.
(8) As in part (7), we have $((x, 0)(0, y))_{0, j}=(x, 0)_{0, j}(0, y)_{0, j}$, for any $j \in \mathbb{N}$. Since $x \geq 1$, we have $(x, 0)_{0, j}=1$, and so $((x, 0)(0, y))_{0, j}$ is $(0, y)_{0, j}$. Thus we have $(x, 0)(0, y)=(0, y)$.
(9) As in part (8), we have $(\infty(0, y))_{0, j}=(0, y)_{0, j}$, and so $\infty(0, y)=\infty$.
(10) By part (1), the element $(1,0)$ is the multiplicative identity; so the result is clear for $x=1$. If $x \geq 2$, then (8.3.1) and (3.4) imply

$$
d(\infty(x, 0)) \succcurlyeq \psi(\infty) d(x, 0)=\infty(x-1,0)
$$

which evaluates to $\infty$ by induction. Therefore we have $d(\infty(x, 0))=\infty$ and hence $\infty(x, 0)=\infty$.

Proposition 8.4 Consider the canonical map $W(\mathbb{B}) \rightarrow W_{(p)}(\mathbb{B})$, in the coordinates given by the bijections (4.8.1) and (8.2.2). If $x \leq 1$ (or rather $x \leq 2$ for $p=2$ ), then an element $(x, y) \in \mathbb{N}^{2}$ is sent to $(x, z)$, where $z$ is the smallest element of $\mathbb{N}$ such that $p^{z}>y$. All other elements are sent to $\infty$.

Proof It is clear that $\infty$ is sent to $\infty$. For elements of the form $(x, 0)$, the result follows from part (1) of (8.3). Now consider an element of the form $(0, y) \in W(\mathbb{B})$. Then $(0, y)\left(\psi_{p^{j}}\right)=0$ if and only if $p^{j}>y$. Therefore the image $a \in W_{(p)}(\mathbb{B})$ of $(0, y)$ is $(0, z)$, where $z$ is the smallest element of $\mathbb{N}$ such that $p^{z}>y$.

For an element of the form $(x, y) \in W(\mathbb{B})$, we have $(x, y)=(x, 0)+(0, y)$. So the result follows from the previous cases and parts (4) and (5) of (8.3).

### 8.5 Remarks

The set of $p$-typical symmetric functions of length $k$ is defined to be the sub- $\mathbb{N}$ algebra of $\Lambda_{\mathbb{N}}$ generated by the subset $\left\{d_{i, j} \mid i+j \leq k\right\}$, and the functor it represents is defined to be $W_{(p), k}$, the $p$-typical Witt functor of length $k$. Then we have the following analogue of (8.2.1):
$W_{(p), k}(\mathbb{B})=\left\{\left(a_{i, j}\right) \in \mathbb{B}^{T_{k}} \mid a_{i, j}=0 \Leftrightarrow a_{i, j+1}=a_{i+1, j}=0 \quad\right.$ for $\left.i+j \leq k-1\right\}$,
where $T_{k}=\left\{(i, j) \in \mathbb{N}^{2} \mid i+j \leq k\right\}$. But there is no analogue of (8.2.2). Indeed, the triangular array

$$
\left(a_{i, j}\right)=\left(\begin{array}{lll}
0 & & \\
1 & 1 & \\
1 & 1 & 0
\end{array}\right)
$$

has two blocks of zeros yet defines an element of $W_{(p), 2}(\mathbb{B})$.
This also shows that the map $W_{(p), 4}(\mathbb{B}) \rightarrow W_{(p), 3}(\mathbb{B})$ is not surjective, because the element above does not lift to $W_{(p), 4}(\mathbb{B})$. In fact, as one easily checks, this element lifts to $W_{(p), 3}\left(\mathbb{R}_{+}\right)$, and so the map $W_{(p), 4}\left(\mathbb{R}_{+}\right) \rightarrow W_{(p), 3}\left(\mathbb{R}_{+}\right)$is not surjective either.

## 9 Complements: countability of some Witt vector algebras

If $A$ is uncountable, then both $W(A)$ and $W^{S c h}(A)$ are uncountable because the Teichmüller map $A \rightarrow W(A)$ is injective. If $A$ is countable, we have the following partial result when it admits an embedding into $\mathbb{R}_{+}$:
Theorem 9.1 Let A be a countable sub- $\mathbb{N}$-algebra of $\mathbb{R}_{+}$.
(1) If $A$ is not dense, then $W(A)$ and $W^{\text {Sch }}(A)$ are countable.
(2) If $A$ is dense and of the form $B \cap \mathbb{R}_{+}$for some subring $B \subseteq \mathbb{R}$, then $W(A)$ and $W^{\text {Sch }}(A)$ are uncountable.

Proof (1) As in (7.5), the series associated to any element of $W^{\text {Sch }}(A)$ is of the form $f(t) / g(t)$, where $f(t)$ and $g(t)$ are polynomials with coefficients in the smallest subfield $K \subseteq \mathbb{R}$ containing $A$. Since $K$ is countable, there are only countably many such series. Further, $W(A)$ is countable because it is a subset of $W^{\operatorname{Sch}}(A)$.
(2) For the same reason, here it is enough to show $W(A)$ is uncountable. First observe that a Witt vector $y \in W(\mathbb{R})$ lies in $W\left(\mathbb{R}_{+}\right)$if and only if, for all $k \geq 0$, there exists a Witt vector $x \in W\left(\mathbb{R}_{+}\right)$such that $\sigma(y) \equiv \sigma(x) \bmod t^{k+1}$. Indeed, to have $y \in W\left(\mathbb{R}_{+}\right)$it is enough to have $y\left(m_{\lambda}\right) \geq 0$ for all partitions $\lambda$. But each such $m_{\lambda}$ is a polynomial in finitely many elementary symmetric functions $e_{1}, \ldots, e_{k}$. Choosing $x \in W\left(\mathbb{R}_{+}\right)$such that $\sigma(y) \equiv \sigma(x) \bmod t^{k+1}$, we have $y\left(m_{\lambda}\right)=x\left(m_{\lambda}\right) \geq 0$.

Combining this with the equality $W\left(B \cap \mathbb{R}_{+}\right)=W(B) \cap W\left(\mathbb{R}_{+}\right)$, we see that it is enough to show there are uncountably many series in $1+t B[[t]]$ that can be approximated arbitrarily well in the $t$-adic topology by series of the form $\sigma(x)$, where
$x \in W\left(\mathbb{R}_{+}\right)$. In fact, it will be convenient to consider approximations by a slightly restricted class of series. So first let $W^{\prime}(\mathbb{R})$ denote the subset of $1+t \mathbb{R}[[t]]$ consisting of series of the form $\prod_{j}\left(1+b_{j} t\right)$, for some reals $b_{1}>b_{2}>\cdots>0$. Then by (6.3), we have $W^{\prime}(\mathbb{R}) \subseteq \sigma\left(W\left(\mathbb{R}_{+}\right)\right)$. Second, for any integer $k \geq 0$, let $W_{k}^{\prime}(B)$ denote the intersection of $1+t B[[t]] /\left(t^{k+1}\right)$ and the image of the truncation map

$$
W^{\prime}(\mathbb{R}) \rightarrow 1+t \mathbb{R}[[t]] /\left(t^{k+1}\right), \quad f(t) \mapsto f(t) \bmod t^{k+1}
$$

Then the inverse limit $\lim _{k} W_{k}^{\prime}(B)$, viewed as a subset of $1+t \mathbb{R}[[t]]$, consists of series with coefficients in $B$ that can be approximated arbitrarily well in the $t$-adic topology by elements of $W^{\prime}(\mathbb{R})$. By the remarks above, it is contained in $\sigma\left(W\left(B \cap \mathbb{R}_{+}\right)\right)$.

Therefore it is enough to show that $\lim _{k} W_{k}^{\prime}(B)$ is uncountable. Since $W_{0}^{\prime}(B)$ is nonempty, it is enough to show that the fibers of the truncation maps $W_{k}^{\prime}(B) \rightarrow$ $W_{k-1}^{\prime}(B)$ have infinitely many elements (or even two). So consider an element $\bar{f}(t) \in$ $W_{k-1}^{\prime}(B)$, and choose a lift $f(t)=\prod_{j}\left(1+b_{j} t\right) \in W^{\prime}(\mathbb{R})$. Then for all sufficiently small $\varepsilon \in \mathbb{R}$, the polynomial $\varepsilon t^{k}+\prod_{j=1}^{k}\left(1+b_{j} t\right)$ is still of the form $\prod_{j=1}^{k}\left(1+b_{j}^{\prime} t\right)$, where $b_{1}^{\prime}>\cdots>b_{k}^{\prime}>b_{k+1}$. For such $\varepsilon$, write

$$
f_{\varepsilon}(t)=\prod_{j=1}^{k}\left(1+b_{j}^{\prime} t\right) \prod_{j=k+1}^{\infty}\left(1+b_{j} t\right) \in W^{\prime}(\mathbb{R})
$$

Then we have

$$
\begin{equation*}
f_{\varepsilon}(t)=f(t)+\varepsilon t^{k} \bmod t^{k+1} \tag{9.1.1}
\end{equation*}
$$

Therefore the series $f_{\varepsilon}(t)$ are distinct for distinct $\varepsilon$, and they reduce to $\bar{f}(t)$ modulo $t^{k}$. Further, we have $f_{\varepsilon}(t) \in W_{k}^{\prime}(B)$ if (and only if) the $k$ th coefficient of $f_{\varepsilon}(t)$ lies in $B$. Since $A$ is dense in $\mathbb{R}_{+}$and since $B$ contains $\pm A$, we know that $B$ is dense in $\mathbb{R}$. Therefore Eq. (9.1.1) guarantees that there are infinitely many sufficiently small $\varepsilon \in \mathbb{R}$ such that the $k$ th coefficient of $f_{\varepsilon}(t)$ lies in $B$. In other words, the fiber of the map $W_{k}^{\prime}(B) \rightarrow W_{k-1}^{\prime}(B)$ over the arbitrary element $\bar{f}(t)$ is infinite.

### 9.2 Examples

Since one can embed $\mathbb{N}\left[x_{1}, x_{2}, \ldots\right]$ as a nondense sub- $\mathbb{N}$-algebra of $\mathbb{R}_{+}$, we see that $W(A)$ and $W^{\text {Sch }}(A)$ are countable for any sub- $\mathbb{N}$-algebra $A$ of $\mathbb{N}\left[x_{1}, x_{2}, \ldots\right]$. In particular, we have countability for $A=\mathbb{N}[\alpha]$, where $\alpha$ is any transcendental number. This holds even when $|\alpha|<1$, in which case $A$ is dense in $\mathbb{R}_{+}$. So the condition $A=B \cap \mathbb{R}_{+}$in (9.1)(2) cannot be removed. Similarly, $W(A)$ and $W^{\text {Sch }}(A)$ are countable for the dense subalgebra $A=\mathbb{N}[2-\sqrt{2}] \subseteq \mathbb{R}_{+}$, because $2-\sqrt{2}$ is conjugate to $2+\sqrt{2}$, which is greater than 1 .

Examples where $W(A)$ and $W^{\text {Sch }}(A)$ are uncountable include $A=\mathbb{N}[1 / m]$ for any $m=2,3, \ldots$ and $A=\mathbb{N}[\sqrt{2}-1]=\mathbb{Z}[\sqrt{2}] \cap \mathbb{R}_{+}$.

Finally, we emphasize that there are dense sub- $\mathbb{N}$-algebras $\mathbb{R}_{+}$to which (9.1)(2) does not apply—for example $\mathbb{N}[\sqrt{3}-1]$. At the time of this writing, we do not know whether either $W(A)$ or $W^{\text {Sch }}(A)$ is countable for any such algebras $A$.

Theorem 9.3 Let $n$ be a nonnegative integer, and write $A=\mathbb{N} /(n=n+1)$.
(1) $W(A)$ is countable if and only if $n \leq 2$.
(2) $W^{\text {Sch }}(A)$ is countable if and only if $n \leq 2$.

Note the contrast between the uncountability for $n \geq 3$ and the countability of $W(\mathbb{N})$ and $W^{\text {Sch }}(\mathbb{N})$, given by (6.3) and (1.2).

Proof (1) When $n=0$, the result is clear, and when $n=1$, it follows immediately from (4.8).

Now consider the case $n=2$. For any $a \in W(A)$, write

$$
\operatorname{pker}(a)=\left\{\lambda \in P \mid a\left(m_{\lambda}\right)=0\right\} .
$$

A basic result we will use more than once is the following. Let $\lambda$ be a partition, and write $d_{\lambda}(a)=\sum_{v \notin \operatorname{pker}(a)} b_{\lambda \lambda}^{\nu}$, where $m_{\lambda}^{2}=\sum_{v \in P} b_{\lambda \lambda}^{\nu} m_{\nu}$; then we have

$$
\begin{equation*}
d_{\lambda}(a) \geq 2 \Longrightarrow a\left(m_{\lambda}\right)=2 \tag{9.3.1}
\end{equation*}
$$

Indeed, since all elements in $A$ are multiplicatively idempotent, we have

$$
a\left(m_{\lambda}\right)=a\left(m_{\lambda}\right)^{2}=a\left(m_{\lambda}^{2}\right)=\sum_{\nu} b_{\lambda \lambda}^{\nu} a\left(m_{\nu}\right) \geq d_{\lambda}(a) \geq 2
$$

and hence $a\left(m_{\lambda}\right)=2$.
Now consider the map $\varphi: W(A) \rightarrow W(\mathbb{B})$ induced by the unique $\mathbb{N}$-algebra map $A \rightarrow \mathbb{B}$. Since $W(\mathbb{B})$ is countable, it is enough to show that for each $z \in W(\mathbb{B})$, there are only finitely many $a \in W(A)$ such that $\varphi(a)=z$.

First consider the case $z=\infty$. Then for any nonzero partition $\lambda$, the symmetric function $m_{\lambda}^{2}$ has at least two monomial constituents, by (4.5). Further, neither constituent is in $\operatorname{pker}(a)$ since $\operatorname{pker}(a)=\emptyset$. Therefore we have $a\left(m_{\lambda}\right)=2$, by (9.3.1). Since we also have $a\left(m_{0}\right)=a(1)=1$, there is at most one Witt vector $a$ such that $\varphi(a)=\infty$.

It remains to consider the case $z=(x, y)$, where $x, y \in \mathbb{N}$. It is enough to show that any $a \in \varphi^{-1}(z)$ is determined by its values $a\left(m_{u^{x}}\right)$ for $u \leq 2 y+1$, because there are only finitely many possibilities for such values. We will do this by saying explicitly how these values determine each $a\left(m_{\lambda}\right)$; we will do this by considering a nested sequence of cases in $\lambda$.

If $\lambda_{x+1} \geq y+1$, then we have $\lambda \in \operatorname{pker}(a)$ and hence $a\left(m_{\lambda}\right)$ is zero. In particular, the value of $a\left(m_{\lambda}\right)$ is determined. Therefore we may assume $\lambda_{x+1} \leq y$. First suppose $\lambda_{x+1}>0$. Clearly $m_{\lambda}^{2}$ has the monomial constituent

$$
v=\left(2 \lambda_{1}, \ldots, 2 \lambda_{x}, \lambda_{x+1}, \lambda_{x+1}, \ldots\right) .
$$

Since $\lambda_{x+1}>0$, if we view $\nu$ as a vector, we can write it as the sum of two distinct permutations of $\lambda$. In other words, the $m_{\nu}$ term in $m_{\lambda}^{2}$ is a cross term, and so its coefficient $b_{\lambda \lambda}^{v}$ is at least 2 . Further we have $v_{x+1}=\lambda_{x+1} \leq y$, and so $v \notin \operatorname{pker}(a)$. Therefore we have $d_{\nu}(a) \geq 2$, and hence $a\left(m_{\lambda}\right)$ can only be 2, by (9.3.1). In particular, its value is determined when $\lambda_{x+1}>0$.

Therefore we may assume $\lambda_{x+1}=0$. Now suppose there exists $i$ with $1 \leq i \leq x-1$ such that $\lambda_{i}>\lambda_{i+1}$. Then $m_{\lambda}^{2}$ has two monomial constituents $\nu$ :

$$
\left(2 \lambda_{1}, \ldots, 2 \lambda_{x}\right), \quad\left(2 \lambda_{1}, \ldots, 2 \lambda_{i-1}, \lambda_{i}+\lambda_{i+1}, \lambda_{i+1}+\lambda_{i}, 2 \lambda_{i+2}, \ldots, 2 \lambda_{x}\right) .
$$

Since $\lambda_{i}>\lambda_{i+1}$, the two differ at position $i$ and are hence distinct. Also, both have length at most $x$. So neither lies in $I_{(x, y)}=\operatorname{pker}(a)$. Thus we have $d_{\lambda}(a) \geq 2$, and again by (9.3.1) the value of $a\left(m_{\lambda}\right)$ is determined.

Therefore we may assume there is no such $i$, and hence that $\lambda=u^{x}$, for some $u$. For any integers $p, q \geq y+1$, the only constituent of $m_{p^{x}} m_{q^{x}}$ which does not belong to $I_{(x, y)}$ is $(p+q)^{x}$, and this constituent appears with coefficient 1. This implies $f\left(m_{p^{x}}\right) f\left(m_{q^{x}}\right)=f\left(m_{(p+q)^{x}}\right)$, and so by induction the values $f\left(m_{u^{x}}\right)$ for $u=y+1, \ldots, 2 y+1$ determine $f\left(m_{u^{x}}\right)$ for any $u \geq 2 y+2$. This finishes the proof when $n=2$.

Finally, consider the case $n \geq 3$.
Let $U$ be any subset of the set of all nonzero partitions. Clearly, there are uncountably many choices for this $U$. To show $W(A)$ is uncountable, we will construct an element of $W(A)$ corresponding to $U$ and show that all such elements of $W(A)$ are distinct.

Define an $\mathbb{N}$-linear map $f: \Lambda_{\mathbb{N}} \rightarrow A$ by

$$
f\left(m_{\lambda}\right)= \begin{cases}1 & \text { if } \lambda=0 \\ n-1 & \text { if } \lambda \in U \\ n & \text { otherwise }\end{cases}
$$

for every partition $\lambda$. Clearly, the subset $U$ can be uniquely reconstructed from this map $f$, since $n-1 \neq n$ in $A$. We will now show that $f$ is an element of $W(A)$. To prove this, we must verify that $f$ is an $\mathbb{N}$-algebra homomorphism. Since $f(1)=f\left(m_{0}\right)=1$, it is enough to prove that $f\left(m_{\lambda} m_{\mu}\right)=f\left(m_{\lambda}\right) f\left(m_{\mu}\right)$ for any partitions $\lambda$ and $\mu$.

If $\lambda=0$ or $\mu=0$, this is obvious, since $m_{0}=1$. So we may assume $\lambda$ and $\mu$ are nonzero. Therefore we have

$$
f\left(m_{\lambda}\right) f\left(m_{\mu}\right) \geq(n-1)(n-1) \geq n
$$

since $n \geq 3$. Because every element of $A$ which is at least $n$ must be equal to $n$, this yields $f\left(m_{\lambda}\right) f\left(m_{\mu}\right)=n$.

On the other hand, (4.5)(3) shows that $m_{\lambda} m_{\mu}$ has at least two distinct nonzero constituents. That is, there exist two distinct nonzero partitions $\alpha$ and $\beta$ such that $m_{\lambda} m_{\mu} \geq m_{\alpha}+m_{\beta}$. Therefore we have

$$
f\left(m_{\lambda} m_{\mu}\right) \geq f\left(m_{\alpha}+m_{\beta}\right)=f\left(m_{\alpha}\right)+f\left(m_{\beta}\right) \geq(n-1)+(n-1) \geq n
$$

since $n \geq 3$ again. As above, this yields $f\left(m_{\lambda} m_{\mu}\right)=n$ and hence

$$
f\left(m_{\lambda} m_{\mu}\right)=n=f\left(m_{\lambda}\right) f\left(m_{\mu}\right) .
$$

Thus $f$ is an element of $W(A)$.
Therefore $W(A)$ has at least as many elements as there are subsets of nonzero partitions. Hence, it has uncountably many.
(2) For $n \neq 2$, the arguments above also work in the Schur basis, instead of the monomial one. One simply appeals to (5.12) instead of (4.8), and (5.6)(3) instead of (4.5)(3). When $n=2$, the modifications we need to make are more elaborate. So we shall not record them here.

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