# A SEMIGROUP-THEORETICAL VIEW OF DIRECT-SUM DECOMPOSITIONS AND ASSOCIATED COMBINATORIAL PROBLEMS 

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#### Abstract

Let $R$ be a ring and let $\mathcal{C}$ be a small class of right $R$-modules which is closed under finite direct sums, direct summands, and isomorphisms. Let $\mathcal{V}(\mathcal{C})$ denote a set of representatives of isomorphism classes in $\mathcal{C}$ and, for any module $M$ in $\mathcal{C}$, let $[M]$ denote the unique element in $\mathcal{V}(\mathcal{C})$ isomorphic to $M$. Then $\mathcal{V}(\mathcal{C})$ is a reduced commutative semigroup with operation defined by $[M]+[N]=[M \oplus N]$, and this semigroup carries all information about direct-sum decompositions of modules in $\mathcal{C}$. This semigroup-theoretical point of view has been prevalent in the theory of direct-sum decompositions since it was shown that if $\operatorname{End}_{R}(M)$ is semilocal for all $M \in \mathcal{C}$, then $\mathcal{V}(\mathcal{C})$ is a Krull monoid. Suppose that the monoid $\mathcal{V}(\mathcal{C})$ is Krull with a finitely generated class group (for example, when $\mathcal{C}$ is the class of finitely generated torsion-free modules and $R$ is a one-dimensional reduced Noetherian local ring). In this case we study the arithmetic of $\mathcal{V}(\mathcal{C})$ using new methods from zero-sum theory. Furthermore, based on module-theoretic work of Lam, Levy, Robson, and others we study the algebraic and arithmetic structure of the monoid $\mathcal{V}(\mathcal{C})$ for certain classes of modules over Prüfer rings and hereditary Noetherian prime rings.


## 1. Introduction

The overarching goal of this manuscript is to study direct-sum decompositions of modules into indecomposable modules. Let $R$ be a ring and let $M$ be a right $R$-module. If $M$ is Noetherian or Artinian, then a well-known and simple argument shows that $M$ is a finite direct sum of indecomposable right $R$-modules. If, for example, $M$ is either both Noetherian and Artinian or if $R$ is a principal ideal domain and $M$ is finitely generated, then such a direct-sum decomposition is unique; that is, the Krull-Remak-Schmidt-Azumaya property (KRSA, for short) holds. For a simple example of non-unique direct-sum decomposition, consider a commutative domain $R$ with distinct non-principal maximal ideals $\mathfrak{m}_{1}$ and $\mathfrak{m}_{2}$. Then the epimorphism $\mathfrak{m}_{1} \oplus \mathfrak{m}_{2} \rightarrow R$, defined by $\left(m_{1}, m_{2}\right) \mapsto m_{1}+m_{2}$, gives rise to a split short exact sequence $0 \rightarrow \mathfrak{m}_{1} \cap \mathfrak{m}_{2} \rightarrow \mathfrak{m}_{1} \oplus \mathfrak{m}_{2} \rightarrow R \rightarrow 0$. Hence $\mathfrak{m}_{1} \oplus \mathfrak{m}_{2} \cong R \oplus\left(\mathfrak{m}_{1} \cap \mathfrak{m}_{2}\right)$. Since $\mathfrak{m}_{1}$ and $\mathfrak{m}_{2}$ are not principal and since $R$ is a domain, each $\mathfrak{m}_{i}$ is indecomposable as an $R$-module and is not isomorphic to $R$. Since the pioneering work of Krull, Remak, Schmidt, and Azumaya, direct-sum decompositions have been a classic topic in module theory. We refer the reader to [10] for an overview of the celebrated Krull-Remak-Schmidt-Azumaya Theorem and related topics in direct-sum theory.

The work of Facchini, Herbera, and Wiegand [14, 15, 17] introduced a new semigroup-theoretical approach to the study of direct-sum decompositions of modules when KRSA fails to hold. Let $\mathcal{C}$ be a small class of right $R$-modules which is closed under finite direct sums, direct summands, and isomorphisms. Let $\mathcal{V}(\mathcal{C})$ denote a set of representatives of isomorphism classes in $\mathcal{C}$ and, for any module $M$ in $\mathcal{C}$, let $[M]$ denote the unique element in $\mathcal{V}(\mathcal{C})$ isomorphic to $M$. Then $\mathcal{V}(\mathcal{C})$ is a reduced commutative semigroup with operation defined by $[M]+[N]=[M \oplus N]$, and this semigroup carries all information about direct-sum decompositions

[^0]of modules in $\mathcal{C}$. In particular, $[M]$ is an irreducible element of the semigroup $\mathcal{V}(\mathcal{C})$ if and only if $M$ is an indecomposable module, and direct-sum decompositions of modules in $\mathcal{C}$ are unique (equivalently, KRSA holds) if and only if $\mathcal{V}(\mathcal{C})$ is a free abelian monoid. Suppose $\mathcal{C}$ is a class of $R$-modules as just defined and that KRSA fails. Then arithmetical questions of the following type naturally arise.

Q1: If $M_{1} \oplus M_{2}=N_{1} \oplus \cdots \oplus N_{l}$, where $M_{1}, M_{2}, N_{1}, \ldots, N_{l}$ are indecomposable modules, does there exist an upper bound for $l$ depending only on $\mathcal{C}$ ?
Q2: Suppose an indecomposable module $M$ is isomorphic to a direct summand of $N_{1} \oplus \cdots \oplus N_{l}$ for indecomposable modules $N_{1}, \ldots, N_{l}$. Is there an upper bound (depending only on $\mathcal{C}$ ) for the number $|I|$ such that $M$ is already isomorphic to a direct summand of $\oplus_{i \in I} N_{i}$ ?

We propose the following overall strategy to tackle these and other arithmetical questions regarding nonunique direct-sum decompositions of modules.
A. Use module-theoretic results in order to describe the algebraic structure of the semigroup $\mathcal{V}(\mathcal{C})$.
B. Use factorization theory to study the arithmetic structure of the semigroup $\mathcal{V}(\mathcal{C})$.

This strategy is relatively new, but has been used in several recent papers for certain classes of modules (see, for example [2]). In the present paper we pursue this strategy for three classes of finitely generated modules: torsion-free modules over one-dimensional reduced commutative Noetherian local rings (Section 4), modules over Prüfer rings (Section 5), and right-modules over hereditary Noetherian prime rings (Section 6).

First, suppose that $\mathcal{C}$ is a class of $R$-modules such that the endomorphism ring $\operatorname{End}_{R}(M)$ is semilocal for each $R$-module $M$. In this case, Facchini proved [9, Theorem 3.4] that $\mathcal{V}(\mathcal{C})$ is a reduced Krull monoid. Earlier results in this direction can be found in [14, 15, 53]). A reduced Krull monoid $H$ is uniquely determined by its characteristic $\left(G,\left(m_{g}\right)_{g \in G}\right)$ where $G$ is the class group of $H$ and, for each $g \in G, m_{g}$ is the cardinality of the set of prime divisors lying in the class $g \in G$ (see Section 2). Many arithmetical problems depend only on the set $G_{p}=\left\{g \in G: m_{g}>0\right\}$ of classes containing prime divisors and, for simplicity, we restrict our discussion to this case. Therefore, in order to determine the structure of $\mathcal{V}(\mathcal{C})$, it is required to determine the characteristic of $\mathcal{V}(\mathcal{C})$, or at least the tuple $\left(G, G_{P}\right)$. In general, this is an herculean task. Indeed, even for specific classes $\mathcal{C}$ of modules where it is known that $\operatorname{End}_{R}(M)$ is semilocal for each $M$ in $\mathcal{C}$, except for in very special situations, we have limited information about $\left(G, G_{P}\right)$. For an overview what is known for certain classes of finitely generated modules over certain one- and two-dimensional Noetherian local rings, see [2]. Nevertheless, suppose we are in a situation where we are able to determine the tuple $\left(G, G_{P}\right)$. Then, by a well-known transfer homomorphism (see Proposition 2.3), arithmetical problems in $\mathcal{V}(\mathcal{C})$ can be studied in the (combinatorial) Krull monoid $\mathcal{B}\left(G_{P}\right)$, the monoid of zero-sum sequences over $G_{P}$. Therefore, in this setting, the study of uniqueness and non-uniqueness of direct-sum decompositions can be reduced to zerosum theory, a flourishing subfield of combinatorial and additive number theory. Except for occasional work (see [1] for early contributions), the focus of (arithmetical) zero-sum theory has been restricted to the case where $G$ is a finite abelian group, whereas class groups $G$ stemming from monoids of modules $\mathcal{V}(\mathcal{C})$ are often infinite. In Section 3 we study zero-sum theory over finitely generated free abelian groups using new methods from matroid theory (goal B). The focus will be on the study of the Davenport constant which can often be used to provide bounds on important factorization-theoretic invariants. Specifically, upper and lower bounds on the Davenport constant are given in Theorem 3.13 . In Section 4 these results will be applied to Krull monoids. In particular, module-theoretic results will be used in order to describe $\mathcal{V}(\mathcal{C})$ (goal $\mathbf{A}$ ) in such a way that we can apply the results from Section 3 (see Corollaries 4.6 and 4.7).

Apart from cases where $\mathcal{V}(\mathcal{C})$ is Krull and some trivial cases (say, where KRSA holds whence $\mathcal{V}(\mathcal{C})$ is free abelian), the structure of $\mathcal{V}(\mathcal{C})$ has not been studied from a semigroup-theoretical point of view. Our goal in the present paper is to take this approach for certain classes of finitely generated modules. The most simple case is the classical Theorem of Steinitz which determines the structure of direct-sum decompositions of finitely generated modules over Dedekind domains. This result has found generalizations into many
directions. We consider two generalizations in the setting of certain classes of modules over Prüfer rings and over hereditary Noetherian prime rings.

In Section 5 we study the class of finitely generated projective modules over a class of Prüfer rings. Based on work of Feng and Lam ([18]) we show that $\mathcal{V}\left(\mathcal{C}_{\text {proj }}\right)$ is a finitely primary monoid (and goal $\mathbf{A}$ is achieved). Since the arithmetic of finitely primary monoids has been well-studied, there are well-known answers to questions about the arithmetic of $\mathcal{V}(\mathcal{C})$ (goal $\mathbf{B}$ ). These results are summarized in Theorems 5.1 and 5.3

In Section 6 we study the class of finitely generated projective modules over hereditary Noetherian prime rings (which generalize non-commutative Dedekind prime rings). Deep module theoretic work by Levy and Robson (41) allows an algebraic characterization of the associated monoids of stable isomorphism classes of modules (goal A). We first introduce monoids of this type in an abstract setting and study their arithmetic (Propositions 6.1 and 6.4). We then apply these results to monoids of modules (goal B) over hereditary Noetherian prime rings in Theorem 6.5.

In both the setting of Sections 5 and of Section 6, the respective monoids of modules are seen to be halffactorial (all direct-sum decompositions of a given module have the same length). However, these modules behave very differently with respect to finer arithmetical invariants including the $\omega$-invariants and the tame degrees (see Theorems 5.1, 5.3, and 6.5).

Section 2 is preparatory in nature. There we gather together arithmetical concepts from factorization theory as well as required material on (generalized) Krull monoids and monoids of modules. An even more detailed description of these concepts and their relevance to module theory can be found in [2]. We also refer the reader to the monograph [28] for more information on factorization theory, and to [3] for a friendly introduction to the interplay of factorization theory and module theory.

## 2. Preliminaries

We denote by $\mathbb{N}$ the set of positive integers and set $\mathbb{N}_{0}=\mathbb{N} \cup\{0\}$. For real numbers $a, b \in \mathbb{R}$ we set $[a, b]=\{x \in \mathbb{Z}: a \leq x \leq b\}$. For a subset $L \subset \mathbb{Z}$, we denote by $\Delta(L) \subset \mathbb{N}$ the set of distances of $L$. This is the set $\{l-k: k<l \in L$ and $L \cap[k, l]=\{k, l\}\}$. Let $G$ be an additive abelian group and let $A, B \subset G$ be subsets. Then $A+B=\{a+b: a \in A, b \in B\}$ denotes the sumset of $A$ and $B,-A=\{-a: a \in A\}$ is the negative of $A, g+A=\{g\}+A$ for $g \in G$, and $\langle A\rangle \subset G$ denotes the subgroup of $G$ generated by $A$.

A family $\left(e_{i}\right)_{i \in I}$ of elements of $G$ is said to be independent if $e_{i} \neq 0$ for all $i \in I$ and, for every family $\left(m_{i}\right)_{i \in I} \in \mathbb{Z}^{(I)}$,

$$
\sum_{i \in I} m_{i} e_{i}=0 \quad \text { implies } \quad m_{i} e_{i}=0 \quad \text { for all } \quad i \in I .
$$

An independent family $\left(e_{i}\right)_{i \in I}$ is called a basis for $G$ if $G=\bigoplus_{i \in I}\left\langle e_{i}\right\rangle$. If $G$ is torsion-free, then $\mathrm{r}(G)=$ $\operatorname{dim}_{\mathbb{Q}}\left(G \otimes_{\mathbb{Z}} \mathbb{Q}\right)$ denotes the rank of $G$ and

$$
G_{|I|}^{+}=\left\{\sum_{i \in I} \epsilon_{i} e_{i}: \epsilon_{i} \in\{0,1\} \text { for all } i \in I, \text { but not all equal to zero }\right\} \subset \oplus_{i \in I}\left\langle e_{i}\right\rangle \cong \mathbb{Z}^{(I)}
$$

denotes the set of non-zero vertices of the hypercube in $\oplus_{i \in I}\left\langle e_{i}\right\rangle$. This definition clearly depends on the chosen basis, but throughout we refer to $G_{|I|}^{+}$only after a basis has been fixed.

By a monoid we always mean a commutative semigroup with identity which satisfies the cancellation law. Thus, if $R$ is a commutative ring and $R^{\bullet}$ its set of regular elements, then $R^{\bullet}$ is a multiplicative monoid. Let $H$ be a (multiplicatively written) monoid. We denote by $\mathrm{q}(H)$ a quotient group of $H$, by $H^{\times}$the group of invertible elements of $H$, and by $\widehat{H}$ the complete integral closure of $H$ with $H \subset \widehat{H} \subset \mathrm{q}(H)$. An element $u \in H$ is called an atom if $u \notin H^{\times}$and $u=a b$ with $a, b \in H$ implies that either $a \in H^{\times}$or $b \in H^{\times}$. The set of all atoms of $H$ is denoted by $\mathcal{A}(H)$. We say that a monoid $H$ with identity 1 is reduced if $H^{\times}=\{1\}$, and we denote by $H_{\text {red }}=\left\{a H^{\times}: a \in H\right\}$ the associated reduced monoid of $H$.

Free abelian monoids and groups. A monoid $F$ is free abelian with basis $P \subset F$ if every $a \in F$ has a unique representation of the form

$$
a=\prod_{p \in P} p^{\mathrm{v}_{p}(a)} \quad \text { with } \quad \mathrm{v}_{p}(a) \in \mathbb{N}_{0} \text { and } \mathrm{v}_{p}(a)=0 \text { for almost all } p \in P
$$

In this case, we set $F=\mathcal{F}(P)$. The isomorphism $a \mapsto\left(\mathrm{v}_{p}(a)\right)_{p \in P}$ from the multiplicative monoid $\mathcal{F}(P)$ to the additive monoid $\left(\mathbb{N}_{0}^{(P)},+\right)$ induces an isomorphism from the quotient group $\mathrm{q}(F)$ to $\mathbb{Z}^{(P)}$. We denote by $\mathcal{F}_{\text {rat }}(P)$ the multiplicative monoid isomorphic to $\left(\mathbb{Q}_{\geq 0}^{(P)},+\right)$ and we tacitly assume that $\mathcal{F}(P) \subset \mathcal{F}_{\text {rat }}(P)$. The quotient group of $\mathcal{F}_{\text {rat }}(P)$ is isomorphic to $\left(\mathbb{Q}^{(P)},+\right)$ and an element $a \in \mathrm{q}\left(\mathcal{F}_{\text {rat }}(P)\right)$ will be written in the form $a=\prod_{p \in P} p^{\vee_{p}(a)} \quad$ with $\quad \mathrm{v}_{p}(a) \in \mathbb{Q}$ and $\mathrm{v}_{p}(a)=0$ for almost all $p \in P$. We call

$$
|a|=\sum_{p \in P}\left|\mathrm{v}_{p}(a)\right| \in \mathbb{Q} \geq 0 \quad \text { the length of } a \text { and } \operatorname{supp}(a)=\left\{p \in P: \mathrm{v}_{p}(a) \neq 0\right\} \subset P \quad \text { the support of } a
$$

Clearly there are inclusions $\mathbb{N}_{0}^{(P)} \subset \mathbb{Z}^{(P)} \subset \mathbb{Q}^{(P)} \subset \mathbb{Q}^{P}$. Elements of $\mathbb{Q}^{P}$ will (usually) be written as $\boldsymbol{x}$ and, if $\boldsymbol{x} \in \mathbb{Q}^{P}$, then we tacitly assume that $\boldsymbol{x}=\left(x_{p}\right)_{p \in P}$ and $\mathrm{v}_{p}(\boldsymbol{x})=x_{p}$ for all $p \in P$. For $p \in P$, let $\boldsymbol{e}_{p} \in \mathbb{N}_{0}^{(P)}$ denote the standard vector with $e_{p, q}=1$ if $p=q$ and $e_{p, q}=0$ for all $q \in P \backslash\{p\}$. Then $\left(\boldsymbol{e}_{p}\right)_{p \in P}$ is the standard basis of $\mathbb{Z}^{(P)}$.

Factorizations and sets of lengths. Let $H$ be a monoid. The free abelian monoid $\mathrm{Z}(H)=\mathcal{F}\left(\mathcal{A}\left(H_{\text {red }}\right)\right)$ is called the factorization monoid of $H$ and the unique homomorphism

$$
\pi: \mathrm{Z}(H) \rightarrow H_{\text {red }} \quad \text { satisfying } \quad \pi(u)=u \quad \text { for each } \quad u \in \mathcal{A}\left(H_{\mathrm{red}}\right)
$$

is called the factorization homomorphism of $H$. For $a \in H$,

$$
\begin{aligned}
& \mathrm{Z}_{H}(a)=\mathrm{Z}(a)=\pi^{-1}\left(a H^{\times}\right) \subset \mathrm{Z}(H) \text { is the set of factorizations of } a, \\
& \mathrm{~L}_{H}(a)=\mathrm{L}(a)=\{|z|: z \in \mathrm{Z}(a)\} \subset \mathbb{N}_{0} \text { is the set of lengths of } a, \text { and } \\
& \mathcal{L}(H)=\{\mathrm{L}(a): a \in H\} \text { is the system of sets of lengths of } H .
\end{aligned}
$$

We say that $H$ is atomic if $\mathrm{Z}(a) \neq \emptyset$ for each $a \in H$, that $H$ is an FF-monoid if $\mathrm{Z}(a)$ is finite and nonempty for each $a \in H$, and that $H$ is factorial if $|\mathrm{Z}(a)|=1$ for each $a \in H$. For the remainder of this section we assume that $H$ is atomic.

Among the most well-studied invariants in factorization theory are those that describe the structure of sets of lengths of elements in $H$. Let $k \in \mathbb{N}$. If $H \neq H^{\times}$, then

$$
\mathcal{U}_{k}(H)=\bigcup_{k \in L \in \mathcal{L}(H)} L
$$

is the union of sets of lengths containing $k$. If $H^{\times}=H$, we set $\mathcal{U}_{k}(H)=\{k\}$. In either case we define $\rho_{k}(H)=\sup \mathcal{U}_{k}(H)$ and $\lambda_{k}(H)=\min \mathcal{U}_{k}(H)$. Clearly, $\mathcal{U}_{1}(H)=\{1\}$ and, for each $k \in \mathbb{N}, k \in \mathcal{U}_{k}(H)$. In particular, $k+l \in \mathcal{U}_{k}(H)+\mathcal{U}_{l}(H) \subset \mathcal{U}_{k+l}(H)$ for each $k, l \in \mathbb{N}$. With $\Delta(L)$ the set of distances of a length set $L$,

$$
\Delta(H)=\bigcup_{L \in \mathcal{L}(H)} \Delta(L)
$$

denotes the set of distances of $H$. By definition, $\rho_{k}(H)=k$ for all $k \in \mathbb{N}$ if and only if $\mathcal{U}_{k}(H)=\{k\}$ for all $k \in \mathbb{N}$ if and only if $\Delta(H)=\emptyset$. In this case, $H$ is said to be half-factorial. If $\Delta(H)=\{d\}$ for some $d \in \mathbb{N}$, then for all $L \in \mathcal{L}(H)$ and for all $k \in \mathbb{N}, L$ and $\mathcal{U}_{k}(H)$ are arithmetical progressions with difference $d$.

Let $M \in \mathbb{N}_{0}, d \in \mathbb{N}$, and $\{0, d\} \subset \mathcal{D} \subset[0, d]$. A subset $L \subset \mathbb{Z}$ is called an almost arithmetical multiprogression (AAMP) with difference $d$, period $\mathcal{D}$, and bound $M$ if

$$
L=y+\left(L^{\prime} \cup L^{*} \cup L^{\prime \prime}\right) \subset y+\mathcal{D}+d \mathbb{Z} \text { where }
$$

- $L^{*}$ is finite and nonempty with $\min L^{*}=0$ and $L^{*}=(\mathcal{D}+d \mathbb{Z}) \cap\left[0, \max L^{*}\right]$,
- $L^{\prime} \subset[-M,-1]$,
- $L^{\prime \prime} \subset \max L^{*}+[1, M]$, and
- $y \in \mathbb{Z}$.

Note that every AAMP is a finite non-empty subset of $\mathbb{Z}$ and that an AAMP with period $\{0, d\}$ and bound $M=0$ is a (usual) arithmetical progression with difference $d$.
Distance between factorizations and catenary degrees. Let $z, z^{\prime} \in \mathrm{Z}(H)$. Then we can write

$$
z=u_{1} \cdot \ldots \cdot u_{l} v_{1} \cdot \ldots \cdot v_{m} \quad \text { and } \quad z^{\prime}=u_{1} \cdot \ldots \cdot u_{l} w_{1} \cdot \ldots \cdot w_{n}
$$

where $l, m, n \in \mathbb{N}_{0}$ and $u_{1}, \ldots, u_{l}, v_{1}, \ldots, v_{m}, w_{1}, \ldots, w_{n} \in \mathcal{A}\left(H_{\text {red }}\right)$ are such that

$$
\left\{v_{1}, \ldots, v_{m}\right\} \cap\left\{w_{1}, \ldots, w_{n}\right\}=\emptyset .
$$

Then $\operatorname{gcd}\left(z, z^{\prime}\right)=u_{1} \cdot \ldots \cdot u_{l}$ and we call

$$
\mathrm{d}\left(z, z^{\prime}\right)=\max \{m, n\}=\max \left\{\left|z \operatorname{gcd}\left(z, z^{\prime}\right)^{-1}\right|,\left|z^{\prime} \operatorname{gcd}\left(z, z^{\prime}\right)^{-1}\right|\right\} \in \mathbb{N}_{0}
$$

the distance between $z$ and $z^{\prime}$. The catenary degree $\mathrm{c}(a)$ of an element $a \in H$ is the smallest $N \in \mathbb{N}_{0} \cup\{\infty\}$ such that, for any two factorizations $z$ and $z^{\prime}$ of $a$, there exists a finite sequence $z=z_{0}, z_{1}, \ldots, z_{k}=z^{\prime}$ of factorizations of $a$ such that $\mathrm{d}\left(z_{i-1}, z_{i}\right) \leq N$ for all $i \in[1, k]$. We denote by

$$
\mathrm{c}(H)=\sup \{\mathrm{c}(a): a \in H\} \in \mathbb{N}_{0} \cup\{\infty\}
$$

the catenary degree of $H$. By definition, $|\mathrm{Z}(a)|=1$ (i.e., $a$ has unique factorization in $H$ ) if and only if $\mathrm{c}(a)=0$, and thus $H$ is factorial if and only if $\mathrm{c}(H)=0$. Suppose now that $H$ is not factorial. Then it is an easy consequence of the definitions that $2+\sup \Delta(H) \leq \mathrm{c}(H)$. In particular, if $\mathrm{c}(H)=2$, then $\Delta(H)=\emptyset$, and if $\mathrm{c}(H)=3$, then $\Delta(H)=\{1\}$.
The $\omega$-invariants and tame degrees. We now recall the $\omega$-invariants as well as local and global tame degrees. These are well-studied invariants in the factorization theory of rings and semigroups. They have also been considered in module-theoretic situations in terms of the so-called semi-exchange property (see [6]).

For an element $a$ in an atomic monoid $H$, let $\omega(H, a)$ denote the smallest $N \in \mathbb{N}_{0} \cup\{\infty\}$ having the following property:

For any multiple $b$ of $a$ and any factorization $b=v_{1} \cdot \ldots \cdot v_{n}$ of $b$, there exists a subset $\Omega \subset[1, n]$ with $|\Omega| \leq N$ such that

$$
a \mid \prod_{\nu \in \Omega} v_{\nu}
$$

We then define

$$
\omega(H)=\sup \{\omega(H, u): u \in \mathcal{A}(H)\} \in \mathbb{N}_{0} \cup\{\infty\}
$$

It is then clear that an atom $u \in H$ is prime if and only if $\omega(H, u)=1$ and thus $H$ is factorial if and only if $\omega(H) \leq 1$. If $H$ satisfies the ascending chain condition (ACC) on divisorial ideals (in particular, if $H$ is Krull or $H=R^{\bullet}$ where $R$ is a Noetherian domain), then $\omega(H, u)<\infty$ for all $u \in \mathcal{A}(H)$ [29, Theorem 4.2].

Roughly speaking, for an element $u$ in an atomic monoid $H$, the tame degree $\mathrm{t}(H, u)$ is the maximum of $\omega(H, u)$ and the factorization lengths of $u^{-1} \prod_{\nu \in \Omega} v_{\nu}$ in $H$. For an atom $u \in H_{\text {red }}$, the local tame degree $\mathrm{t}(H, u)$ is the smallest $N \in \mathbb{N}_{0} \cup\{\infty\}$ with the following property:

For any multiple $a \in H_{\text {red }}$ of $u$ and any factorization $z=v_{1} \cdot \ldots \cdot v_{n} \in \mathbf{Z}(a)$ of $a$ which does not contain $u$, there is a short subproduct which is a multiple of $u$, say $v_{1} \cdot \ldots \cdot v_{m}$, and a refactorization of this subproduct which contains $u$, say $v_{1} \cdot \ldots \cdot v_{m}=u u_{2} \cdot \ldots \cdot u_{\ell}$, such that $\max \{\ell, m\} \leq N$.
In particular, the local tame degree $\mathrm{t}(H, u)$ measures the distance between any factorization of a multiple $a$ of $u$ and a factorization of $a$ which contains $u$. We denote by

$$
\mathrm{t}(H)=\sup \left\{\mathrm{t}(H, u): u \in \mathcal{A}\left(H_{\mathrm{red}}\right)\right\} \in \mathbb{N}_{0} \cup\{\infty\}
$$

the (global) tame degree of $H$, and we say that $H$ is tame if $\mathrm{t}(H)<\infty$. If $u$ is prime, then $\mathrm{t}(H, u)=0$ and thus $H$ is factorial if and only if $\mathrm{t}(H)=0$. In order to describe the relationship between the $\omega$-invariants and the tame degree, we observe that

$$
\begin{gathered}
\omega(H, a)=\sup \left\{k \in \mathbb{N}_{0} \cup\{\infty\}:\right. \\
\quad b=u_{1} \cdot \ldots \cdot u_{k} \in a H \text { with } k \in \mathbb{N}_{0}, u_{1}, \ldots, u_{k} \in \mathcal{A}(H), \\
\\
\text { and } \left.u \nmid u_{i}^{-1} b \text { for all } i \in[1, k]\right\}
\end{gathered}
$$

and we define the $\tau$-invariant as

$$
\begin{align*}
\tau(H, a)=\sup \left\{\min \mathrm{L}\left(a^{-1} b\right):\right. & b=u_{1} \cdot \ldots \cdot u_{k} \in a H \text { with } k \in \mathbb{N}_{0}, u_{1}, \ldots, u_{k} \in \mathcal{A}(H), \\
& \text { and } \left.a \nmid u_{i}^{-1} b \text { for all } i \in[1, k]\right\} . \tag{1}
\end{align*}
$$

For each non-prime $u \in \mathcal{A}(H), \mathrm{t}\left(H, u H^{\times}\right)=\max \{\omega(H, u), \tau(H, u)+1\}$ ([29, Theorem 3.6]). Suppose that $H$ is half-factorial and $\mathrm{L}(a)=\{l\}$ with $l \in \mathbb{N}_{0}$. Then $\mathrm{L}\left(a^{-1} b\right)=\{k-l\}$ and hence $\omega(H, a)=\tau(H, a)+l$. In particular, if $u \in \mathcal{A}(H)$, then $\omega(H, u)=\tau(H, u)+1$.

We recall that if $H_{\text {red }}$ is finitely generated, then $H$ is tame ([28, Theorem 3.1.4]). Tame monoids will be studied in Proposition 2.3. Theorem 5.3, and Theorem 6.5. In Section 4 we will characterize when the monoid $\mathcal{V}((\mathcal{C})$ is tame for $\mathcal{C}$ a class of finitely generated modules over a commutative Noetherian local ring (Proposition 4.1). Further examples of tame monoids can be found in [30, 39]. We now gather together several arithmetical finiteness properties of tame monoids.

## Proposition 2.1 (Arithmetic of tame monoids).

Let $H$ be a tame monoid.

1. If $H$ is not factorial, then $2+\sup \Delta(H) \leq \mathrm{c}(H) \leq \omega(H) \leq \mathrm{t}(H) \leq \omega(H)^{2}<\infty$.
2. There is a constant $M \in \mathbb{N}_{0}$ such that every set of lengths $L \in \mathcal{L}(H)$ is an AAMP with difference $d \in \Delta(H)$ and bound $M$.
3. There is a constant $M \in \mathbb{N}_{0}$ such that, for every $k \geq 2$, the union $\mathcal{U}_{k}(H)$ of sets of lengths is an AAMP with period $\{0, \min \Delta(H)\}$ and bound $M$.

Proof. Statement 1 follows easily from [28, Theorem 1.6.3] and [30, Section 3]. Statement 2 follows from 30, Theorem 5.1] and statement 3 follows from [24, Theorems 4.2 and 3.5].

Transfer homomorphisms. A monoid homomorphism $\theta: H \rightarrow B$ is called a transfer homomorphism if the following properties are satisfied:
(T 1) $B=\theta(H) B^{\times}$and $\theta^{-1}\left(B^{\times}\right)=H^{\times}$.
(T 2) If $u \in H, b, c \in B$, and $\theta(u)=b c$, then there exist $v, w \in H$ such that $u=v w, \theta(v) \simeq b$, and $\theta(w) \simeq c$.
Transfer homomorphisms are a central tool in Factorization Theory and allow one to lift arithmetical results from a (simpler) monoid $B$ to the monoid $H$ (of original interest). These homomorphisms will be used throughout this manuscript (see, in particular, Propositions 2.3. 6.1. and 6.4. Each transfer homomorphism naturally gives rise to a homomorphism $\bar{\theta}$ of factorization monoids that extends $\theta$. Let $\theta: H \rightarrow B$ be a transfer homomorphism of atomic monoids. Then $\theta$ induces a homomorphism $\bar{\theta}: \mathbf{Z}(H) \rightarrow \mathbf{Z}(B)$ satisfying $\bar{\theta}\left(u H^{\times}\right)=\theta(u) B^{\times}$for all $u \in \mathcal{A}(H)$.

We now recall how various factorization-theoretic invariants are preserved by transfer homomorphisms. Let $\theta: H \rightarrow B$ and $\bar{\theta}: \mathbf{Z}(H) \rightarrow \mathbf{Z}(B)$ be as above. For $a \in H$, we denote by $c(a, \theta)$ the smallest $N \in \mathbb{N}_{0} \cup\{\infty\}$ with the following property:

If $z, z^{\prime} \in \mathrm{Z}_{H}(a)$ and $\bar{\theta}(z)=\bar{\theta}\left(z^{\prime}\right)$, then there exist $k \in \mathbb{N}_{0}$ and factorizations $z=z_{0}, \ldots, z_{k}=z^{\prime} \in$ $\mathrm{Z}_{H}(a)$ such that $\bar{\theta}\left(z_{i}\right)=\bar{\theta}(z)$ and $\mathrm{d}\left(z_{i-1}, z_{i}\right) \leq N$ for all $i \in[1, k]$; that is, $z$ and $z^{\prime}$ can be concatenated by an $N$-chain in the fiber $\mathrm{Z}_{H}(a) \cap \bar{\theta}^{-1}(\bar{\theta}(z))$.

Now

$$
\mathrm{c}(H, \theta)=\sup \{\mathrm{c}(a, \theta): a \in H\} \in \mathbb{N}_{0} \cup\{\infty\}
$$

denotes the catenary degree in the fibres of $\theta$. The next lemma summarizes what will be needed in the sequel (a proof can be found in [28, Section 3.2]).

Lemma 2.2. Let $\theta: H \rightarrow B$ be a transfer homomorphism.

1. $H$ is atomic if and only if $B$ is atomic.
2. For all $a \in H, a$ is an atom of $H$ if and only if $\theta(a)$ is an atom of $B$.
3. Suppose that $H$ is atomic.
(a) For all $a \in H, \mathrm{~L}_{H}(a)=\mathrm{L}_{B}(\theta(a))$. In particular,
(i) $\mathcal{L}(H)=\mathcal{L}(B)$ and $\Delta(H)=\Delta(B)$, and
(ii) $\mathcal{U}_{k}(H)=\mathcal{U}_{k}(B), \rho_{k}(H)=\rho_{k}(B)$, and $\lambda_{k}(H)=\lambda_{k}(B)$ for each $k \in \mathbb{N}$.
(b) For all $a \in H, \mathrm{c}_{B}(\theta(a)) \leq \mathrm{c}_{H}(a) \leq \max \left\{\mathrm{c}_{B}(\theta(a)), \mathrm{c}(a, \theta)\right\}$. In particular,

$$
\mathrm{c}(B) \leq \mathrm{c}(H) \leq \max \{\mathrm{c}(B), \mathrm{c}(H, \theta)\} .
$$

Generalized Krull monoids. Let $H$ and $D$ be monoids. A monoid homomorphism $\varphi: H \rightarrow D$ is called

- a divisor homomorphism if $\varphi(a) \mid \varphi(b)$ implies that $a \mid b$ for all $a, b \in H$.
- a divisor theory (for $H$ ) if $D=\mathcal{F}(P)$ for some set $P \subset D, \varphi$ is a divisor homomorphism, and, for every $a \in \mathcal{F}(P)$, there exists a finite nonempty subset $X \subset H$ such that $a=\operatorname{gcd}(\varphi(X))$.
- cofinal if for every $\alpha \in D$ there exists $a \in H$ such that $\alpha \mid \varphi(a)$.

If $H \subset D$ is a submonoid, then $H \subset D$ is said to be saturated (or cofinal) if the embedding $H \hookrightarrow D$ is a divisor homomorphism (or is cofinal). The monoid $H$ is called a

- rational generalized Krull monoid if there is a divisor homomorphism $\varphi: H \rightarrow \mathcal{F}_{\text {rat }}(P)$ for some set $P$.
- Krull monoid if there is a divisor homomorphism $\varphi: H \rightarrow \mathcal{F}(P)$ for some set $P$.

We note that every Krull monoid is a rational generalized Krull monoid and that a monoid $H$ is a (rational generalized) Krull monoid if and only if $H_{\text {red }}$ is a (rational generalized) Krull monoid. Generalized Krull monoids and domains have been studied in [32, Section 5], 37, Chapter 22], and 4]. Specifically, [4, Proposition 2] guarantees that the definition of rational generalized Krull monoids we have given above coincides with the usual one.

We note that a monoid is Krull if and only if it is completely integrally closed and $v$-noetherian and that every Krull monoid is an FF-monoid and has a divisor theory. Let $H$ be a Krull monoid and let $\varphi: H \rightarrow D=\mathcal{F}(P)$ be a cofinal divisor homomorphism. We call $\mathcal{C}(\varphi)=\mathrm{q}(D) / \mathrm{q}(\varphi(H))$ the class group of $\varphi$ and use additive notation for this group. For $a \in \mathrm{q}(D)$, we denote by $[a]=[a]_{\varphi}=a \mathrm{q}(\varphi(H)) \in \mathrm{q}(D) / \mathrm{q}(\varphi(H))$ the class containing $a$. Since $\varphi$ is a cofinal divisor homomorphism, $\mathcal{C}(\varphi)=\{[a]: a \in D\}$ and $\varphi(H)=\{a \in$ $D:[a]=[1]\}$. The set

$$
G_{P}=\{[p]=p \mathbf{q}(\varphi(H)): p \in P\} \subset \mathcal{C}(\varphi)
$$

of classes containing prime divisors plays a crucial role in arithmetic computations (see Proposition 2.3) and we have $\left[G_{P}\right]=\mathcal{C}(\varphi)$. If $\varphi$ is a divisor theory, then $\varphi$ and the class group $\mathcal{C}(\varphi)$ are unique up to isomorphism. Since $\mathcal{C}(\varphi)$ depends only on $H$, we denote it by $\mathcal{C}(H)$ and call it the class group of $H$. Moreover, a reduced Krull monoid $H$ with divisor theory $H \hookrightarrow \mathcal{F}(P)$ is uniquely determined up to isomorphism by its characteristic $\left(G,\left(m_{g}\right)_{g \in G}\right)$ where $G$ is an abelian group together with an isomorphism $\Phi: G \rightarrow \mathcal{C}(H)$ and where $\left(m_{g}\right)_{g \in G}$ is a family of cardinal numbers $m_{g}=|P \cap \Phi(g)|$ (see [28, Theorem 2.5.4]). We consider the characteristic of a monoid of modules $\mathcal{V}(\mathcal{C})$ that is Krull in both Sections 4 and 6 .

A domain $R$ is a rational generalized Krull domain if and only if $R^{\bullet}$ is a rational generalized Krull monoid (for recent work on these kinds of domains, see [42, 47). A $v$-Marot ring (and, in particular, a domain) $R$ is
a Krull ring if and only if its multiplicative monoid of regular elements $R^{\bullet}$ is a Krull monoid ([33, Theorem 3.5]), and we set $\mathcal{C}(R)=\mathcal{C}\left(R^{\bullet}\right)$. If $R$ is a Dedekind domain, then $\mathcal{C}(R)$ is the Picard group of $R$.

We now introduce a Krull monoid of combinatorial flavor, the monoid $\mathcal{B}\left(G_{0}\right)$ of zero-sum sequences over a subset $G_{0}$ of an abelian group $G$. As previously mentioned, this monoid will play a crucial role in arithmetical investigations of general Krull monoids. Section 3 provides a detailed study of $\mathcal{B}\left(G_{0}\right)$ in case of subsets $G_{0}$ of finitely generated free abelian groups. Let $G$ be an additive abelian group, $G_{0} \subset G$ a subset, and $S=g_{1} \cdot \ldots \cdot g_{l} \in \mathcal{F}\left(G_{0}\right)$. We call $\sigma(S)=g_{1}+\cdots+g_{l} \in G$ the sum of $S$, and we define, for $k \in \mathbb{N}$,

$$
\begin{aligned}
\Sigma(S) & =\left\{\sum_{i \in I} g_{i}: \emptyset \neq I \subset[1, \ell]\right\} \subset G \\
\Sigma_{k}(S) & =\left\{\sum_{i \in I} g_{i}: I \subset[1, \ell],|I|=k\right\} \subset G \quad \text { and } \quad \Sigma_{\leq k}(S)=\bigcup_{\nu \in[1, k]} \Sigma_{\nu}(S)
\end{aligned}
$$

If $\left(e_{1}, \ldots, e_{r}\right)$ is a basis of $G$, then the set $G_{r}^{+}$of nonzero vertices of the hypercube satisfies

$$
G_{r}^{+}=\Sigma\left(e_{1} \cdot \ldots \cdot e_{r}\right)
$$

This set will be thoroughly studied in Sections 3 and 4 .
For a map $\varphi: G \rightarrow G^{\prime}$ between two abelian groups $G$ and $G^{\prime}$, we define $\varphi(S)=\varphi\left(g_{1}\right) \cdot \ldots \cdot \varphi\left(g_{\ell}\right)$. Also, for $S \in \mathcal{F}\left(G_{0}\right)$, we set $-S=\left(-g_{1}\right) \cdot \ldots \cdot\left(-g_{\ell}\right)$. Clearly

$$
\mathcal{B}\left(G_{0}\right)=\left\{S \in \mathcal{F}\left(G_{0}\right): \sigma(S)=0\right\} \subset \mathcal{F}\left(G_{0}\right)
$$

is a submonoid of $\mathcal{F}\left(G_{0}\right)$. Moreover, since the inclusion $\mathcal{B}\left(G_{0}\right) \hookrightarrow \mathcal{F}\left(G_{0}\right)$ is a divisor homomorphism, $\mathcal{B}\left(G_{0}\right)$ is a Krull monoid. For arithmetical invariants $*(\cdot)$, as defined previously, we write (as it is usual) $*\left(G_{0}\right)$ instead of $*\left(\mathcal{B}\left(G_{0}\right)\right)$. In particular, $\mathcal{A}\left(G_{0}\right)$ denotes the set of atoms of $\mathcal{B}\left(G_{0}\right), \Delta\left(G_{0}\right)$ denotes the set of distances of $\mathcal{B}\left(G_{0}\right)$, and so forth. Note that the atoms of $\mathcal{B}\left(G_{0}\right)$ are precisely the minimal zero-sum sequences over $G_{0}$ those zero-sum sequences having no proper subsequence that is also a zero-sum sequence - and we denote by

$$
\mathrm{D}\left(G_{0}\right)=\sup \left\{|U|: U \in \mathcal{A}\left(G_{0}\right)\right\} \in \mathbb{N}_{0} \cup\{\infty\}
$$

the Davenport constant of $G_{0}$, a central invariant in zero-sum theory ( $25,51,31,36$ ). The following lemma highlights the close connection between the arithmetic of a general Krull monoid and the arithmetic of the associated monoid of zero-sum sequences. A proof can be found in [28, Theorems 3.4.2 and 3.4.10].

Proposition 2.3. Let $H$ be a Krull monoid, $\varphi: H \rightarrow D=\mathcal{F}(P)$ a cofinal divisor homomorphism, $G=\mathcal{C}(\varphi)$ its class group, and $G_{P} \subset G$ the set of classes containing prime divisors. Let $\widetilde{\boldsymbol{\beta}}: D \rightarrow \mathcal{F}\left(G_{P}\right)$ denote the unique homomorphism defined by $\widetilde{\boldsymbol{\beta}}(p)=[p]$ for all $p \in P$.

1. The homomorphism $\boldsymbol{\beta}=\widetilde{\boldsymbol{\beta}} \circ \varphi: H \rightarrow \mathcal{B}\left(G_{P}\right)$ is a transfer homomorphism. Moreover, $\mathrm{c}(H, \boldsymbol{\beta}) \leq 2, \mathrm{c}\left(G_{P}\right) \leq \mathrm{c}(H) \leq \max \left\{\mathrm{c}\left(G_{P}\right), 2\right\}$, and $\mathrm{c}(H) \leq \mathrm{D}\left(G_{P}\right)$.
2. If $G_{P}$ is finite, then $\mathcal{A}\left(G_{P}\right)$ is finite and hence $\mathrm{D}\left(G_{P}\right)<\infty$.
3. If $\mathrm{D}\left(G_{P}\right)<\infty$, then both $H$ and $\mathcal{B}\left(G_{P}\right)$ are tame.

Monoids of modules. Let $R$ be a ring and let $\mathcal{C}$ be a small class of right $R$-modules. That is, $\mathcal{C}$ has a set $\mathcal{V}(\mathcal{C})$ of isomorphism class representatives, and for any $M \in \mathrm{Ob}(\mathcal{C})$, we denote by $[M]$ the unique element of $\mathcal{V}(\mathcal{C})$ isomorphic to $M$. In more technical terms we suppose that the full subcategory $\mathcal{C}$ of $\operatorname{Mod}-R$ is skeletally small. Suppose that $\mathcal{C}$ is closed under finite direct sums, direct summands, and isomorphisms. Then $\mathcal{V}(\mathcal{C})$ is a reduced commutative semigroup with operation $[M]+[N]=[M \oplus N]$, and all information about direct-sum decomposition of modules in $\mathcal{C}$ can be studied in terms of factorizations in the semigroup $\mathcal{V}(\mathcal{C})$.

Suppose that $\mathcal{V}(\mathcal{C})$ is a monoid and let $\mathcal{C}^{\prime}$ be a subclass of $\mathcal{C}$ which is closed under isomorphisms. Then $\mathcal{C}^{\prime}$ is closed under finite direct sums, direct summands, and isomorphisms if and only if $\mathcal{V}\left(\mathcal{C}^{\prime}\right) \subset \mathcal{V}(\mathcal{C})$ is a divisor-closed submonoid. For a module $M$ in $\mathcal{C}$ we denote by add $(M)$ the class of $R$-modules that are isomorphic to direct summands of direct sums of finitely many copies of $M$. Then $\mathcal{V}(\operatorname{add}(M))$ is the smallest
divisor-closed submonoid generated by $[M] \in \mathcal{V}(\mathcal{C})$. The class $\mathcal{C}_{\text {proj }}$ of finitely generated projective right $R$-modules is of special importance ( 7 ] and [13, Section 2.2]) and will be considered in Sections 5 and 6 .

Proposition 2.4. Let $R$ be a ring and $\mathcal{C}$ a small class of right $R$-modules which is closed under finite direct sums, direct summands, and isomorphisms.

1. If $\operatorname{End}_{R}(M)$ is semilocal for each $M$ in $\mathcal{C}$, then $\mathcal{V}(\mathcal{C})$ is a Krull monoid.
2. If there exists $M$ in $\mathcal{C}$ such that $\operatorname{End}_{R}(M)$ is semilocal, then $\mathcal{V}(\operatorname{add}(M))$ is a finitely generated Krull monoid. Conversely, if $\mathcal{V}(\mathcal{C})$ is a finitely generated monoid, then there is some $M$ in $\mathcal{C}$ such that $\mathcal{V}(\mathcal{C})=\mathcal{V}(\operatorname{add}(M))$.
3. If $R$ is semilocal, then $\mathcal{V}\left(\mathcal{C}_{\text {proj }}\right)$ is a finitely generated Krull monoid. In addition, if $R$ is commutative, then $\mathcal{V}\left(\mathcal{C}_{\text {proj }}\right)$ is Krull if and only if $\mathcal{V}\left(\mathcal{C}_{\text {proj }}\right)$ is free abelian.

Proof. See [9, Theorem 3.4] for the proof of 1.
We now consider statement 2. Let $N_{1}$ and $N_{2}$ be right $R$-modules and recall that $\operatorname{End}\left(N_{1} \oplus N_{2}\right)$ is semilocal if and only if $\operatorname{End}\left(N_{1}\right)$ and $\operatorname{End}\left(N_{2}\right)$ are semilocal (see the introduction of [16]). Thus $\operatorname{End}(N)$ is semilocal for all $N$ in $\operatorname{add}(M)$ and hence $\mathcal{V}(\operatorname{add}(M))$ is a Krull monoid by 1. Moreover, [8, Corollary 4.11] implies that $\mathcal{V}(\operatorname{add}(M))$ is finitely generated. Conversely, suppose that $\mathcal{V}(\mathcal{C})$ is a finitely generated monoid. Let $M_{1}, \ldots, M_{t}$ be indecomposable right $R$-modules such that $\mathcal{A}(\mathcal{V}(\mathcal{C}))=\left\{\left[M_{1}\right], \ldots,\left[M_{t}\right]\right\}$. Then $M=M_{1} \oplus \cdots \oplus M_{t}$ has the required property.

Since $\mathcal{V}\left(\mathcal{C}_{\text {proj }}\right)=\mathcal{V}\left(\operatorname{add}\left(R_{R}\right)\right)$ and $\operatorname{End}\left(R_{R}\right) \cong R$, the first claim in 3 follows immediately from 2. The second statement follows from [12, Theorem 4.2].

## Remark 2.5.

1. Statement 1 of Proposition 2.4 is a local condition guaranteeing that $\mathcal{V}(\mathcal{C})$ is a Krull monoid and such examples will be considered in Section 4 . It is well-known that the condition " $\operatorname{End}_{R}(M)$ is semilocal" is stronger than the condition that $" \mathcal{V}(\operatorname{add}(M))$ is a Krull monoid". Indeed, if $H$ is any Krull monoid, then all divisor-closed submonoids are Krull as well. However, even if all divisor-closed submonoids which are generated by a single element are Krull monoids, then the monoid $H$ can fail to be Krull (see [48] for counterexamples).
2. In Section 5 we will see that since a finitely primary monoid $H$ is Krull if and only if $H$ is factorial, the monoid of modules $\mathcal{V}\left(\mathcal{C}_{\text {proj }}\right)$ in this setting shares the same property.
3. In Section 6 we will study finitely generated monoids of modules which are not Krull.
4. Every reduced Krull monoid is isomorphic to a monoid of finitely generated projective modules $\mathcal{V}\left(\mathcal{C}_{\text {proj }}\right)$ for which $\operatorname{End}_{R}(M)$ is semilocal for all $M \in \mathcal{C}_{\text {proj }}$ ([17, Theorem 2.1]). For further realization results see, for example [53], [13, Section 9], 40], and [28, Theorem 2.7.14].
5. In 11, Facchini presents an extensive list of classes of modules having semilocal endomorphism rings, and thus a list of monoids of modules $\mathcal{V}(\mathcal{C})$ which are Krull.

## 3. Zero-Sum Theory in finitely generated free abelian groups

In this section we study the Davenport constant $\mathrm{D}\left(G_{0}\right)$ when $G_{0}$ is a subset of a finitely generated free abelian group $G$. Although the results we obtain are interesting in their own right, we have in mind an application to the study of invariants of certain Krull monoids, and in particular certain monoids of modules $\mathcal{V}(\mathcal{C})$ that we study in Section 4.

Let $G$ be an additive abelian group and let $G_{0} \subset G$ be a subset. A sequence over $G_{0}$ will mean a finite sequence of terms from $G_{0}$ which is unordered and where repetition of terms is allowed. Zero-Sum Theory studies such sequences, their sets of subsequence sums, and their structure under extremal conditions (see
[23, 34, 36]). Much of the focus has been on sequences over finite abelian groups, but - motivated by applications in various fields - sequences over infinite abelian groups have recently found more attention (see [26, 50] for examples of recent work).

Our goal in this section is to present, and expand upon, known methods from matroid theory that can be used to find upper bounds for the Davenport constant $\mathrm{D}\left(G_{0}\right)$ for subsets $G_{0}$ of finitely generated free abelian groups. These methods from matroid theory were used by Sturmfels with regards to toric varieties [52]. Such toric varieties are, in turn, related to primitive partition identities which were further studied by Diaconis, Graham, and Sturmfels [5] by combining the matroid-theoretic methods with methods from the geometry of numbers. Although no mention of either zero-sums or the Davenport constant was made in their work, Freeze and Schmid [21, Theorem 6.5] along with Geroldinger and Yuan 35, Section 3] made attempts to use these methods in order to study zero-sum problems over finite abelian groups. We continue this program with a goal of studying zero-sum problems over infinite abelian groups.

In order to make these methods from matroid theory - written using very different notation and terminology than is typical in zero-sum theory - more generally available for the study of zero-sum problems over finitely generated free abelian groups, we develop the basic theory here from scratch. In addition, we modify the arguments so that they extend to the case when $G_{0}$ is non-symmetric, a situation not included in the original formulations from [49] and 52].

Let $G$ be a finitely generated free abelian group, $Q$ a $\mathbb{Q}$-vector space with $G \subset Q$, and $G_{0}$ a subset $G$. The elements $S \in \mathcal{F}_{\text {rat }}\left(G_{0}\right)$ will be called rational sequences over $G_{0}$. If $S=\prod_{g \in G_{0}} g^{v_{g}(S)} \in \mathcal{F}_{\text {rat }}\left(G_{0}\right)$, then $\sigma(S)=\sum_{g \in G_{0}} \vee_{g}(S) g \in Q$ is called the sum of $S$. Then

$$
\mathcal{B}_{\mathrm{rat}}\left(G_{0}\right)=\left\{S \in \mathcal{F}_{\mathrm{rat}}\left(G_{0}\right): \sigma(S)=0\right\} \subset \mathcal{F}_{\mathrm{rat}}\left(G_{0}\right)
$$

is a saturated submonoid of $\mathcal{F}_{\text {rat }}\left(G_{0}\right)$, and hence $\mathcal{B}_{\text {rat }}\left(G_{0}\right)$ is a rational generalized Krull monoid. This monoid is clearly reduced and, since for any nonidentity $B \in \mathcal{B}_{\text {rat }}\left(G_{0}\right)$ we have $B=B^{1 / 2} B^{1 / 2}, \mathcal{B}_{\text {rat }}\left(G_{0}\right)$ has no atoms. Moreover, $\mathcal{B}_{\text {rat }}\left(G_{0}\right)$ is a generalized block monoid as introduced in [27, Example 4.10]. The elements of $\mathcal{B}_{\text {rat }}\left(G_{0}\right)$ will be called rational zero-sum sequences over $G_{0}$. Obviously, for any rational (zerosum) sequence $S$ there exists an integer $n \in \mathbb{N}$ such that $S^{n}$ is an ordinary (zero-sum) sequence. Thus $\mathcal{B}\left(G_{0}\right) \subset \mathcal{B}_{\text {rat }}\left(G_{0}\right)$ and $\mathcal{F}\left(G_{0}\right) \subset \mathcal{F}_{\text {rat }}\left(G_{0}\right)$ are root extensions (see [4, Section 5]). If $G_{0} \subset G_{0}^{\prime}$ are two subsets of $G$, we assume $\mathcal{F}_{\text {rat }}\left(G_{0}\right) \subset \mathcal{F}_{\text {rat }}\left(G_{0}^{\prime}\right)$ and $\mathcal{B}_{\text {rat }}\left(G_{0}\right) \subset \mathcal{B}_{\text {rat }}\left(G_{0}^{\prime}\right)$. Specifically, we make this assumption when we consider $G_{0}^{\prime}=G_{0} \cup-G_{0}$. Given a sequence $S \in \mathcal{F}_{\text {rat }}\left(G_{0}\right)$ and $g \in G_{0}$, we tacitly use terms such as $\mathrm{v}_{-g}(S)$ and $-S$, where we interpret $-g$ and $-S$ as elements of $\mathcal{F}_{\text {rat }}\left(G_{0} \cup-G_{0}\right)$. In particular, $\mathrm{v}_{-g}(S)=0$ if $-g \notin G_{0}$.

We begin with the central construction of a partion of $G_{0}$ and of an epimorphism from $\mathcal{F}_{\text {rat }}\left(G_{0}\right)$ to some $\mathbb{Q}$-vector space. This construction will remain valid throughout Section 3 .

Partition the nonzero elements of $G_{0}$ as $G_{0} \backslash\{0\}=G_{0}^{+} \cup G_{0}^{-}$, where the elements of $G_{0}$ have been distributed so that if $g,-g \in G_{0} \backslash\{0\}$, then $g$ and $-g$ neither both occur in $G_{0}^{+}$nor both occur in $G_{0}^{-}$. Furthermore, choose such a partition so that $G_{0}^{+}$is maximal, that is, such that $-G_{0}^{-} \subset G_{0}^{+}$. There may, of course, be many ways to achieve this, but we choose one such partition and fix it for the remainder of Section 3

Let

$$
\varphi^{\prime}: \mathcal{F}_{\text {rat }}\left(G_{0}\right) \rightarrow \mathcal{F}_{\text {rat }}\left(G_{0}^{+} \backslash\left(-G_{0}^{-}\right)\right) \times \mathrm{q}\left(\mathcal{F}_{\text {rat }}\left(-G_{0}^{-}\right)\right) \subset \mathrm{q}\left(\mathcal{F}_{\text {rat }}\left(G_{0}^{+}\right)\right)
$$

be the unique epimorphism satisfying $\varphi^{\prime}(g)=g$ for each $g \in G_{0}^{+}, \varphi^{\prime}(-g)=g^{-1}$ for each $-g \in G_{0}^{-}$, and $\varphi^{\prime}(0)=1$ provided $0 \in G_{0}$. The arguments of this section are based in the geometry of the sets and thus can be phrased more naturally using vector notation. To translate, we use the canonical isomorphism between $\mathrm{q}\left(\mathcal{F}_{\text {rat }}\left(G_{0}^{+}\right)\right)$and $\mathbb{Q}^{\left(G_{0}^{+}\right)}$which maps $g$ to $\mathrm{e}_{g}$ for each $g \in G_{0}^{+}$and where $\left(\mathrm{e}_{g}\right)_{g \in G_{0}^{+}}$denotes the standard basis of $\mathbb{Q}^{\left(G_{0}^{+}\right)}$. Composing $\varphi^{\prime}$ with this isomorphism, we obtain an epimorphism

$$
\varphi: \mathcal{F}_{\mathrm{rat}}\left(G_{0}\right) \xrightarrow{\varphi^{\prime}} \mathcal{F}_{\mathrm{rat}}\left(G_{0}^{+} \backslash\left(-G_{0}^{-}\right)\right) \times \mathrm{q}\left(\mathcal{F}_{\mathrm{rat}}\left(-G_{0}^{-}\right)\right) \rightarrow \mathbb{Q}_{\geq 0}^{\left(G_{0}^{+} \backslash\left(-G_{0}^{-}\right)\right)} \oplus \mathbb{Q}^{\left(-G_{0}^{-}\right)} \subset \mathbb{Q}^{\left(G_{0}^{+}\right)}
$$

satisfying

$$
\varphi\left(\prod_{g \in G_{0}} g^{\mathrm{v}_{g}(S)}\right)=\sum_{g \in G_{0}^{+}}\left(\mathrm{v}_{g}(S)-\mathrm{v}_{-g}(S)\right) \mathrm{e}_{g}
$$

The sequences $S \in \operatorname{ker}(\varphi) \cap \mathcal{F}\left(G_{0}\right)$ are precisely those zero-sum sequences over $G_{0}$ having a factorization into zero-sum subsequences of length at most 2 , and an arbitrary rational sequence from $\operatorname{ker}(\varphi)$ is nothing more than a rational power of such a sequence. The purpose of this construction is to create a setting where we can first apply methods from linear algebra and the geometry of $\mathbb{Q}$-vector spaces to study $\varphi(S)$, and then to apply the results we obtain to the original sequence $S \in \mathcal{F}_{\text {rat }}\left(G_{0}\right)$.

For each rational sequence $S \in \mathcal{F}_{\text {rat }}\left(G_{0}\right)$ we define the signed support of $S$ as

$$
\operatorname{supp}^{+}(S)=\left\{g \in G_{0} \cup-G_{0}: \mathrm{v}_{g}(S)-\mathrm{v}_{-g}(S) \neq 0\right\}
$$

Observe that for each $S \in \mathcal{F}\left(G_{0}\right)$, we have

$$
\operatorname{supp}^{+}(S)=\left\{g \in G_{0} \cup-G_{0}: \varphi^{\prime}(g) \in \operatorname{supp}\left(\varphi^{\prime}(S)\right)\right\}
$$

and

$$
\operatorname{supp}^{+}(S)=\operatorname{supp}^{+}(-S) \subset \operatorname{supp}(S) \cup-\operatorname{supp}(S) \subset G_{0} \cup-G_{0}
$$

For $S \in \mathcal{F}_{\text {rat }}\left(G_{0}\right)$ we define

$$
\begin{equation*}
R=0^{\mathrm{v}_{0}(S)} \prod_{g \in G_{0}^{+}}(g(-g))^{\min \left\{\mathrm{v}_{g}(S), \mathrm{v}_{-g}(S)\right\}} \in \mathcal{F}_{\mathrm{rat}}\left(G_{0}\right), \tag{2}
\end{equation*}
$$

where $\mathrm{v}_{-g}(S)=0$ whenever $-g \notin G_{0}$. Then $R \mid S$ and we set $S^{\prime}=R^{-1} S \in \mathcal{F}_{\text {rat }}\left(G_{0}\right)$. It is then easily noted that $\operatorname{supp}^{+}(S)=\operatorname{supp}^{+}\left(S^{\prime}\right)$, that $\varphi(S)=\varphi\left(S^{\prime}\right)$, that $\varphi(R)=\mathbf{0} \in \mathbb{Q}^{\left(G_{0}^{+}\right)}$, and for each $g \in G_{0}, g$ and $-g$ are not both contained in $\operatorname{supp}\left(S^{\prime}\right)$ for any $g \in G_{0}$; that is, either $\mathrm{v}_{g}\left(S^{\prime}\right)=0$ or $\mathrm{v}_{-g}\left(S^{\prime}\right)=0$. Note that the latter condition is equivalent to

$$
\operatorname{supp}\left(S^{\prime}\right) \cap-\operatorname{supp}\left(S^{\prime}\right)=\emptyset
$$

In particular, $0 \notin \operatorname{supp}\left(S^{\prime}\right)$. Finally, if $S \in \mathcal{F}\left(G_{0}\right)$ is an ordinary sequence, then $R, S^{\prime} \in \mathcal{F}\left(G_{0}\right)$ are also ordinary sequences. These observations will be used repeatedly in the sequel.

We say that a (rational) zero-sum sequence $S$ is elementary if $\operatorname{supp}^{+}(S)$ is nonempty and minimal; that is, there is no (rational) zero-sum sequence $T$ with $\emptyset \neq \operatorname{supp}^{+}(T) \subsetneq \operatorname{supp}^{+}(S)$. Clearly, a rational zero-sum sequence $S$ is elementary if and only if $S^{n}$ is an elementary zero-sum sequence for some positive integer $n$. An elementary atom is an atom $U \in \mathcal{B}\left(G_{0}\right)$ which is also an elementary zero-sum sequence. We let

$$
\mathrm{D}^{\text {elm }}\left(G_{0}\right)=\sup \left\{|U|: U \in \mathcal{A}\left(G_{0}\right) \text { is elementary }\right\} \in \mathbb{N}_{0} \cup\{\infty\}
$$

denote the supremum of the lengths of elementary atoms over $G_{0}$ (with the convention that $\sup \emptyset=0$ ), and call $\mathrm{D}^{\mathrm{elm}}\left(G_{0}\right)$ the elementary Davenport constant over $G_{0}$.

In our motivating applications the group $G$ is the class group of a Krull monoid, and thus we need each of the previously defined concepts in this general abstract setting. However, it is technically simpler but no restriction to work over the additive group $\mathbb{Z}^{r}$ instead of working over an abstract finitely generated free abelian group. (Note, if $\varphi: G \rightarrow \mathbb{Z}^{r}$ is a group isomorphism and $G_{0} \subset G$, then $\mathrm{D}\left(G_{0}\right)=\mathrm{D}\left(\varphi\left(G_{0}\right)\right)$.) Therefore, for the rest of this section we suppose that

$$
G_{0} \subset G=\mathbb{Z}^{r} \quad \text { where } \quad r \geq 1
$$

Since the case $G_{0} \subset\{0\}$ is trivial, we further assume that the set $G_{0}$ contains a nonzero element of $G$. Moreover, whenever it is convenient, we may assume that $\mathrm{r}\left(\left\langle G_{0}\right\rangle\right)=r$, as otherwise we could replace $\mathbb{Z}^{r}$ with $\left\langle G_{0}\right\rangle \cong \mathbb{Z}^{r}\left(\left\langle G_{0}\right\rangle\right)$. We start with a sequence of basic but important properties regarding elementary zero-sum sequences.

Lemma 3.1. Let $r \geq 1$ and let $G_{0} \subset \mathbb{Z}^{r}$ be a nonempty subset. Let $S, T \in \mathcal{F}_{\text {rat }}\left(G_{0}\right)$ be rational sequences such that $\operatorname{supp}(S) \cap \operatorname{supp}(-S)=\emptyset$ and $\operatorname{supp}(T) \cap \operatorname{supp}(-T)=\emptyset$. Suppose $\operatorname{supp}^{+}(T) \subset \operatorname{supp}^{+}(S)$ and $\operatorname{supp}(T) \cap \operatorname{supp}(S) \neq \emptyset$. Let

$$
\alpha=\min \left\{\mathrm{v}_{g}(S) / \mathrm{v}_{g}(T): g \in \operatorname{supp}(S) \cap \operatorname{supp}(T)\right\} \in \mathbb{Q}_{>0} .
$$

Define $S^{\prime}=(-T)^{\alpha} S \in \mathcal{F}_{\text {rat }}\left(G_{0} \cup-G_{0}\right)$, let

$$
R=\prod_{g \in G_{0}^{+}}(g(-g))^{\min \left\{\mathrm{v}_{g}\left(S^{\prime}\right), \mathrm{v}_{-g}\left(S^{\prime}\right)\right\}} \in \mathcal{F}_{\mathrm{rat}}\left(G_{0} \cup-G_{0}\right),
$$

and set $\tilde{S}^{\prime}:=S^{\prime} R^{-1}$. Then $\tilde{S}^{\prime} \in \mathcal{F}_{\text {rat }}\left(G_{0}\right)$ is a rational sequence over $G_{0}$ with

$$
\begin{equation*}
\operatorname{supp}\left(\tilde{S}^{\prime}\right) \subsetneq \operatorname{supp}(S) \quad \text { and } \quad \operatorname{supp}^{+}\left(\tilde{S}^{\prime}\right) \subsetneq \operatorname{supp}^{+}(S) \tag{3}
\end{equation*}
$$

Proof. By definition, $\tilde{S}^{\prime} \in \mathcal{B}_{\text {rat }}\left(G_{0} \cup-G_{0}\right)$ with $\operatorname{supp}^{+}\left(\tilde{S^{\prime}}\right)=\operatorname{supp}^{+}\left(S^{\prime}\right)$. Thus, by the definition of $S^{\prime}$,

$$
\begin{equation*}
\operatorname{supp}^{+}\left(\tilde{S}^{\prime}\right)=\operatorname{supp}^{+}\left(S^{\prime}\right) \subset \operatorname{supp}^{+}(-T) \cup \operatorname{supp}^{+}(S)=\operatorname{supp}^{+}(T) \cup \operatorname{supp}^{+}(S)=\operatorname{supp}^{+}(S) \tag{4}
\end{equation*}
$$

with the final equality obtained by using the hypothesis that $\operatorname{supp}^{+}(T) \subset \operatorname{supp}^{+}(S)$. By hypothesis, $\operatorname{supp}^{+}(T) \subset$ supp $^{+}(S)$, and thus

$$
\begin{equation*}
-g \in \operatorname{supp}(S) \subset G_{0} \quad \text { for every } g \in \operatorname{supp}(T) \backslash \operatorname{supp}(S) \tag{5}
\end{equation*}
$$

Let $g_{1} \in \operatorname{supp}(S) \cap \operatorname{supp}(T)$ with $\mathrm{v}_{g}(S) / \mathrm{v}_{g}(T)$ minimal in $\mathbb{Q}_{>0}$ where $g \in \operatorname{supp}(S) \cap \operatorname{supp}(T)$. By hypothesis, $g_{1}$ and $-g_{1}$ cannot both be in $\operatorname{supp}(S)$ nor can they both be in $\operatorname{supp}(T)$. Thus, since $g_{1} \in \operatorname{supp}(T) \cap \operatorname{supp}(S)$, we see that $-g_{1} \notin \operatorname{supp}(T)$ and $-g_{1} \notin \operatorname{supp}(S)$. Consequently, by definition of $\alpha$, we have

$$
\mathrm{v}_{-g_{1}}\left(-T^{\alpha}\right)=\mathrm{v}_{g_{1}}\left(T^{\alpha}\right)=\mathrm{v}_{g_{1}}(S) \quad \text { and } \quad \mathrm{v}_{g_{1}}\left(-T^{\alpha}\right)=\alpha \mathrm{v}_{-g_{1}}(T)=0=\mathrm{v}_{-g_{1}}(S)
$$

But $g_{1} \in \operatorname{supp}^{+}(S)$ as $g_{1} \in \operatorname{supp}(S)$ and $-g_{1} \notin \operatorname{supp}(S)$. Thus $g_{1} \notin \operatorname{supp}^{+}\left(S^{\prime}\right)$ and the inclusion in (4) must be strict; that is,

$$
\begin{equation*}
\operatorname{supp}^{+}\left(\tilde{S}^{\prime}\right) \subsetneq \operatorname{supp}^{+}(S) \tag{6}
\end{equation*}
$$

From the definitions of $S^{\prime}$ and $\tilde{S}^{\prime}$, we have

$$
\begin{equation*}
\operatorname{supp}\left(\tilde{S}^{\prime}\right) \subset \operatorname{supp}\left(S^{\prime}\right) \subset \operatorname{supp}(-T) \cup \operatorname{supp}(S) \tag{7}
\end{equation*}
$$

For $g \in \operatorname{supp}(T) \cap \operatorname{supp}(S)$, the definition of $\alpha$ ensures that $\mathrm{v}_{-g}\left(-T^{\alpha}\right) \leq \mathrm{v}_{g}(S)$. Since $g$ and $-g$ cannot both be in $\operatorname{supp}(T)$ nor both be in $\operatorname{supp}(S)$, we also have that $\mathrm{v}_{g}\left(-T^{\alpha}\right)=0=\mathrm{v}_{-g}(S)$. Combining this fact with the definition of $\tilde{S}^{\prime}$ we see that for every $g \in \operatorname{supp}(T) \cap \operatorname{supp}(S),-g \notin \operatorname{supp}\left(\tilde{S}^{\prime}\right)$. Considering this fact along with (5) and (7), it follows that

$$
\begin{equation*}
\operatorname{supp}\left(\tilde{S}^{\prime}\right) \subset \operatorname{supp}(S) \subset G_{0} \tag{8}
\end{equation*}
$$

In particular, we now know that $\tilde{S}^{\prime} \in \mathcal{F}_{\text {rat }}\left(G_{0}\right)$. It remains to show that the inclusion in (8) is strict. But for this fact we need only recall that $g_{1} \in \operatorname{supp}(S)$ and that $g_{1} \notin \operatorname{supp}^{+}\left(\widetilde{S^{\prime}}\right)=\operatorname{supp}\left(\widetilde{S^{\prime}}\right) \cup-\operatorname{supp}\left(\widetilde{S^{\prime}}\right)$. Indeed, this last equality follows immediately since $\operatorname{supp}\left(\widetilde{S^{\prime}}\right)$ contains at most one of $g$ and $-g$ for every $g \in G_{0}$, a fact that follows from the definition of $\operatorname{supp}\left(\widetilde{S^{\prime}}\right)$ and the observations before Lemma 3.1.

Lemma 3.2. Let $r \geq 1$ and let $G_{0} \subset \mathbb{Z}^{r}$ be a nonempty subset. If $S, T \in \mathcal{B}_{\mathrm{rat}}\left(G_{0}\right)$ are both elementary rational zero-sum sequences with $\operatorname{supp}^{+}(S)=\operatorname{supp}^{+}(T)$, then either

$$
S R_{T}^{\alpha}=T^{\alpha} R_{S} \quad \text { or } \quad S\left(-R_{T}\right)^{\alpha}=-T^{\alpha} R_{S} \quad \text { for some positive } \alpha \in \mathbb{Q}_{>0},
$$

where $R_{S} \mid S$ and $R_{T} \mid T$ are the respective maximal length rational zero-sum subsequences of $S$ or $T$ with $\varphi\left(R_{T}\right)=\mathbf{0}$ and $\varphi\left(R_{S}\right)=\mathbf{0}$ (defined in (2)).

In particular, if $S, T \in \mathcal{B}\left(G_{0}\right)$ are elementary zero-sum sequences with common signed support, then there exist relatively prime $m, n \in \mathbb{N}$ such that either

$$
S^{n} R_{T}^{m}=T^{m} R_{S}^{n} \quad \text { or } \quad S^{n}\left(-R_{T}\right)^{m}=(-T)^{m} R_{S}^{n}
$$

where $R_{S} \mid S$ and $R_{T} \mid T$ are the respective maximal length zero-sum subsequences of $S$ or $T$ having a factorization into zero-sum subsequences each of length at most 2.
Proof. Let $S, T \in \mathcal{B}_{\text {rat }}\left(G_{0}\right)$ be elementary rational zero-sum sequences with common signed support supp ${ }^{+}(S)=$ supp $^{+}(T)$. In view of the observations before Lemma 3.1 we may, without loss of generality, assume that $R_{S}$ and $R_{T}$ are trivial (the general case follows easily from this special case). In particular,

$$
\begin{equation*}
\operatorname{supp}(S) \cap \operatorname{supp}(-S)=\emptyset \text { and } \operatorname{supp}(T) \cap \operatorname{supp}(-T)=\emptyset \tag{9}
\end{equation*}
$$

Moreover, since $S$ and $T$ are elementary zero-sum sequences, $\operatorname{supp}^{+}(S)=\operatorname{supp}^{+}(T) \neq \emptyset$.
First suppose that $\operatorname{supp}(S) \cap \operatorname{supp}(T) \neq \emptyset$. Then, since $\operatorname{supp}^{+}(S)=\operatorname{supp}^{+}(T)$, we can apply Lemma 3.1 Let $\alpha \in \mathbb{Q}_{>0}, S^{\prime} \in \mathcal{F}_{\text {rat }}\left(G_{0} \cup-G_{0}\right)$, and $\tilde{S}^{\prime} \in \mathcal{F}_{\text {rat }}\left(G_{0}\right)$ be as in Lemma 3.1. Since $S, T$, and $R$ are each rational zero-sum sequences, it follows from the definition of $\tilde{S}^{\prime}$ that $\tilde{S}^{\prime}$ is also a rational zero-sum sequence. Thus $\tilde{S}^{\prime}$ will, in view of $(3)$, contradict that $S$ is an elementary rational zero-sum sequence unless supp ${ }^{+}\left(\tilde{S}^{\prime}\right)=\emptyset$. However, by definition, supp ${ }^{+}\left(\tilde{S}^{\prime}\right)=\emptyset$ is only possible if $\tilde{S}^{\prime}$ is trivial, in which case $\left(-T^{\alpha}\right) S=S^{\prime}=R$. But this implies, in view of (9) and $\operatorname{supp}^{+}(T)=\operatorname{supp}^{+}(S)$, that

$$
\mathrm{v}_{g}\left(T^{\alpha}\right)=\mathrm{v}_{g}(S) \quad \text { for all } g \in G_{0}
$$

Therefore $S=T^{\alpha}$ as desired.
We now assume that $\operatorname{supp}(S) \cap \operatorname{supp}(T)=\emptyset$. In this case, since $\operatorname{supp}^{+}(S)=\operatorname{supp}^{+}(T)$, it follows from (9) that $\operatorname{supp}(T)=-\operatorname{supp}(S)$. This in turn implies that $\operatorname{supp}(-T) \subset G_{0}$ and hence $-T \in \mathcal{B}_{\text {rat }}\left(G_{0}\right)$. Repeating the arguments of the previous paragraph using $-T$ in place of $T$, we conclude that $S=(-T)^{\alpha}$ as desired. The in particular statement follows easily from the general statement with $\alpha=\frac{m}{n}$, where $m, n \in \mathbb{N}$ are relatively prime.

Lemma 3.3. Let $r \geq 1$ and let $G_{0} \subset \mathbb{Z}^{r}$ be a nonempty subset. If $U, V \in \mathcal{A}\left(G_{0}\right)$ are elementary atoms with $\operatorname{supp}^{+}(U)=\operatorname{supp}^{+}(V)$, then either $U=V$ or $U=-V$.

Proof. Since $U$ and $V$ are both elementary, their signed support must be nontrivial. Moreover, since $U$ and $V$ are atoms, neither $U$ nor $V$ has a nontrivial zero-sum subsequence of length one or two. Therefore, applying Lemma 3.2 to both $U$ and $V$, we find that either $U^{n}=V^{m}$ or $U^{n}=(-V)^{m}$ for relatively prime positive integers $m$ and $n$. Note that if the latter case holds, then $\operatorname{supp}(-V)=\operatorname{supp}(U) \subset G_{0}$. Thus, by replacing $V$ with $-V$ if need be (in which case the hypotheses of the theorem hold for $U$ and $-V$ and, if the conclusion holds for $U$ and $-V$, then it will hold for the original pair $U$ and $V$ as well), we may assume that $U^{n}=V^{m}$ for relatively prime positive integers $m$ and $n$. Therefore $n \mathrm{v}_{g}(U)=m \mathrm{v}_{g}(V)$ for all $g \in G_{0}$. Consequently, since $\operatorname{gcd}(m, n)=1$, it follows that $m \mid \mathrm{v}_{g}(U)$ for each $g \in G_{0}$. But then $U=W^{m}$ is a product decomposition of $U$, where $W=\prod_{g \in G_{0}} g^{v_{g}(U) / m} \in \mathcal{F}\left(G_{0}\right)$. Noting that $0=\sigma(U)=m \sigma(W)$, we conclude that $W$ is a zero-sum sequence. Since it was assumed that $U \in \mathcal{A}\left(G_{0}\right)$ is an atom, $m=1$. A similar argument shows that $n=1$. Now the relation $U^{m}=V^{m}$ gives the desired conclusion $U=V$.

Lemma 3.4. Let $r \geq 1$ and let $G_{0} \subset \mathbb{Z}^{r}$ be a nonempty subset. If $S \in \mathcal{B}\left(G_{0}\right)$ is an elementary zero-sum sequence, then

$$
S=R U^{\ell}
$$

for some $\ell \geq 1$, some elementary atom $U \in \mathcal{A}\left(G_{0}\right)$ with $\operatorname{supp}^{+}(U)=\operatorname{supp}^{+}(S)$, and some zero-sum sequence $R \in \operatorname{ker}(\varphi)$ that has a factorization involving only zero-sum subsequences all of length at most 2.

Proof. Assume for the sake of contradiction that $S \in \mathcal{B}\left(G_{0}\right)$ is an elementary zero-sum sequence that fails to have the desired form and with $|S|$ minimal. By the minimality of $|S|$ and the observations made before Lemma 3.1, it follows that every nontrivial zero-sum subsequence of $S$ has length at least 3 . Therefore $\operatorname{supp}(S) \cap \operatorname{supp}(-S)=\emptyset$. Since the conclusion of the lemma holds for each atom, it must be the case that the chosen $S$ is not an atom. Therefore

$$
S=T_{1} T_{2}
$$

for nontrivial zero-sum subsequences $T_{1}$ and $T_{2}$. Now, since every zero-sum subsequence of $S$ has length at least 3 , the nontrivial zero-sum subsequences $T_{1}$ and $T_{2}$ must each have non-empty signed support. But this contradicts the fact that $S$ is an elementary zero-sum sequence unless we have

$$
\begin{equation*}
\operatorname{supp}^{+}(S)=\operatorname{supp}^{+}\left(T_{1}\right)=\operatorname{supp}^{+}\left(T_{2}\right) \tag{10}
\end{equation*}
$$

Also, $\left|T_{1}\right|<|S|$ and $\left|T_{2}\right|<|S|$ since $T_{1}$ and $T_{2}$ are both nontrivial subsequences of $S$. By (10), each $T_{i}$ is a zero-sum sequence with $\operatorname{supp}^{+}\left(T_{i}\right)=\operatorname{supp}^{+}(S)$ and is consequently an elementary zero-sum sequence. Furthermore, since $\left|T_{1}\right|<|S|$ and $\left|T_{2}\right|<|S|$, the minimality of $|S|$ ensures that both $T_{1}$ and $T_{2}$ have the form stated in the conclusion of the lemma. Thus $T_{1}=U^{m}$ and $T_{2}=V^{n}$ for atoms $U, V \in \mathcal{A}\left(G_{0}\right)$ with

$$
\operatorname{supp}^{+}(U)=\operatorname{supp}^{+}(V)=\operatorname{supp}^{+}\left(T_{1}\right)=\operatorname{supp}^{+}\left(T_{2}\right)=\operatorname{supp}^{+}(S)
$$

and positive integers $m$ and $n$. As $U$ and $V$ are both atoms with $\operatorname{supp}^{+}(U)=\operatorname{supp}^{+}(V)=\operatorname{supp}^{+}(S)$, it follows that $U$ and $V$ are both elementary atoms. Invoking Lemma 3.3 we find that either $U=V$ or $U=-V$.

If $U=-V$, then there must exist $g \in \operatorname{supp}(U)$ with $-g \in \operatorname{supp}(V)$. In this case, since $T_{1}=U^{m}, T_{2}=V^{n}$, and $S=T_{1} T_{2}$, it follows that $g(-g) \mid S$, contrary to conclusion above that $S$ does not have any zero-sum subsequence of length 1 or 2 . Therefore $U \neq-V$, forcing $U=V$. Then $T_{1}=U^{m}, T_{2}=V^{n}=U^{n}$, and $S=T_{1} T_{2}=U^{m+n}$ as desired.

Combining the above results, we are now able to describe the form of elementary zero-sum sequences. We do so in the following proposition whose proof follows immediately from Lemmas 3.3 and 3.4 Essentially, Proposition 3.5 states that if $X$ is the signed support of an elementary zero-sum sequence, then (up to sign) there is a unique atom $U$ with $\operatorname{supp}^{+}(U)=X$, and all other zero-sum sequences having signed support $X$ must have the form $U^{\ell} R$ or $(-U)^{\ell} R$ where $\ell \geq 1$ and $R \in \mathcal{B}\left(G_{0}\right)$ is a zero-sum sequence that has a factorization involving zero-sum subsequences all of length at most 2.

Proposition 3.5. Let $r \geq 1$ and let $G_{0} \subset \mathbb{Z}^{r}$ be a nonempty subset. If $X$ is the signed support of an elementary zero-sum sequence over $G_{0}$, then there exists a unique (up to sign) atom $U \in \mathcal{A}\left(G_{0}\right)$ with $\operatorname{supp}^{+}(U)=$ supp $^{+}(-U)=X$ such that every elementary zero-sum sequence $S$ with signed support $\operatorname{supp}^{+}(S)=X$ has the form

$$
S=R U^{\ell} \quad \text { or } \quad S=R(-U)^{\ell}
$$

where $R \in \operatorname{ker}(\varphi)$ is a zero-sum sequence and $\ell \geq 1$.

Lemma 3.6. Let $r \geq 1$ and let $G_{0} \subset \mathbb{Z}^{r}$ be a nonempty subset. Suppose $S \in \mathcal{B}_{\mathrm{rat}}\left(G_{0}\right)$ is a rational zerosum sequence with $\operatorname{supp}^{+}(S)$ nonempty. Then there exists some elementary atom $U \in \mathcal{A}\left(G_{0}\right)$ such that $\operatorname{supp}(U) \subset \operatorname{supp}(S)$.

Proof. Assume for the sake of contradiction that $S \in \mathcal{B}_{\text {rat }}\left(G_{0}\right)$ is a counterexample with $\operatorname{supp}^{+}(S) \neq \emptyset$ minimal. By removing rational zero-sum subsequences of the form $((-g) g)^{\min \left\{\mathrm{v}_{g}(S), \mathrm{v}_{-g}(S)\right\}}$ and $0^{\mathrm{v}_{0}(S)}$ from $S$ as defined in (2), we may, without loss of generality, assume that $\operatorname{supp}(S) \cap \operatorname{supp}(-S)=\emptyset$. Let $U \in \mathcal{B}\left(G_{0}\right)$ be an elementary zero-sum sequence with

$$
\operatorname{supp}^{+}(U) \subset \operatorname{supp}^{+}(S)
$$

Note that such an elementary zero-sum sequence exists with supp ${ }^{+}(U)=X$ for any minimal nonempty subset $X \subset \operatorname{supp}^{+}(S)$ provided there exists a zero-sum sequence with signed support $X$.

In view of the observations made prior to Lemma 3.1, we may assume that $U$ has no nontrivial zero-sum sequence of length 1 or 2 . Now, if $U=U_{1} \ldots \cdot U_{\ell}$ is a factorization as a product of atoms $U_{i} \in \mathcal{A}\left(G_{0}\right)$, we find that $\operatorname{supp}^{+}\left(U_{i}\right)=\operatorname{supp}^{+}(U)$ for each $i \in[1, l]$, else $U$ is not an elementary zero-sum sequence. Therefore, replacing $U$ by some $U_{i}$ as need be, we may without loss of generality assume that $U \in \mathcal{A}\left(G_{0}\right)$ is an atom.

Since $\operatorname{supp}^{+}(U) \subset \operatorname{supp}^{+}(S), \operatorname{supp}(U)$ and $\operatorname{supp}(S)$ are disjoint only if $-g \in \operatorname{supp}(S) \subset G_{0}$ for every $g \in \operatorname{supp}(U)$. In this case, $-U \in \mathcal{A}\left(G_{0}\right)$ is also an atom over $G_{0}$. Thus, replacing $U$ by $-U$ (in this one scenario), we may assume that $\operatorname{supp}(U) \cap \operatorname{supp}(S) \neq \emptyset$. Since an elementary atom cannot have a nontrivial
zero-sum sequence of length 1 or 2 , we may apply Lemma 3.1 with $T=U$. Now let $\alpha, S^{\prime}, R$ and $\tilde{S}^{\prime}$ be as in Lemma 3.1. Since $S, U$ and $R$ are each zero-sum sequences, it follows from the definition of $\tilde{S}^{\prime}$ that $\tilde{S}^{\prime}$ is also a zero-sum sequence. Therefore

$$
\tilde{S}^{\prime} \in \mathcal{B}_{\mathrm{rat}}\left(G_{0}\right)
$$

an improvement over $\tilde{S}^{\prime} \in \mathcal{F}_{\text {rat }}\left(G_{0}\right)$ given in Lemma 3.1.
If $\operatorname{supp}^{+}\left(\tilde{S}^{\prime}\right)=\operatorname{supp}^{+}\left(S^{\prime}\right)$ is nonempty, then by the strict inclusion $\operatorname{supp}^{+}\left(\tilde{S}^{\prime}\right) \subsetneq \operatorname{supp}^{+}(S)$ in (3), we may apply the induction hypothesis to $\tilde{S}^{\prime} \in \mathcal{B}_{\text {rat }}\left(G_{0}\right)$ and find an elementary atom $V \in \mathcal{A}\left(G_{0}\right)$ such that $\operatorname{supp}(V) \subset \operatorname{supp}\left(\tilde{S}^{\prime}\right)$. The proof is then complete if one considers the other inclusion in (3). Therefore we may assume the alternative, that

$$
\begin{equation*}
\operatorname{supp}^{+}\left(\tilde{S^{\prime}}\right)=\operatorname{supp}^{+}\left(S^{\prime}\right)=\emptyset \tag{11}
\end{equation*}
$$

Recalling that $S^{\prime}=(-U)^{\alpha} S$, that $\operatorname{supp}^{+}(U) \subset \operatorname{supp}^{+}(S)$, and that $\operatorname{supp}(S) \cap \operatorname{supp}(-S)=\operatorname{supp}(U) \cap$ $\operatorname{supp}(-U)=\emptyset$, we see that 11) is possible only if $\operatorname{supp}(U)=\operatorname{supp}(S)$ with $\mathrm{v}_{g}\left(U^{\alpha}\right)=\mathrm{v}_{g}(S)$ for every $g \in \operatorname{supp}(U)=\operatorname{supp}(S)$.

Let $G_{0} \subset \mathbb{Z}^{r} \subset \mathbb{Q}^{r}$ be a finite subset. Note that the chosen partition of $G_{0} \backslash\{0\}$ gives rise to a unique partition of $G_{0} \cup-G_{0} \backslash\{0\}$ with $\left(G_{0} \cup-G_{0}\right)^{+}=G_{0}^{+}$. To this partition we may again associate a map $\mathcal{F}_{\text {rat }}\left(G_{0} \cup-G_{0}\right) \rightarrow \mathbb{Q}^{\left(G_{0}^{+}\right)}$which we also denote by $\varphi$. Then $\varphi: \mathcal{F}_{\text {rat }}\left(G_{0}\right) \rightarrow \mathbb{Q}^{\left(G_{0}^{+}\right)}$is simply the restriction of $\varphi: \mathcal{F}_{\text {rat }}\left(G_{0} \cup-G_{0}\right) \rightarrow \mathbb{Q}^{\left(G_{0}^{+}\right)}$to $\mathcal{F}_{\text {rat }}\left(G_{0}\right)$. By construction, we have

$$
\varphi\left(\mathcal{B}\left(G_{0}\right)\right) \subset \mathbb{Z}^{\left(G_{0}^{+}\right)} \subset \mathbb{Q}^{\left(G_{0}^{+}\right)} .
$$

It is easily checked that $\varphi\left(\mathcal{B}_{\text {rat }}\left(G_{0}\right)\right)$ is an additive monoid closed under multiplication by nonnegative rational numbers. The $\mathbb{Q}$-vector space spanned by $\varphi\left(\mathcal{B}_{\text {rat }}\left(G_{0}\right)\right)$ is then $\varphi\left(\mathcal{B}_{\text {rat }}\left(G_{0} \cup-G_{0}\right)\right)$, which is also the $\mathbb{Q}$-vector space spanned by $\varphi\left(\mathcal{B}\left(G_{0}\right)\right)=\varphi\left(\mathcal{B}_{\text {rat }}\left(G_{0}\right)\right) \cap \mathbb{Z}^{\left(G_{0}^{+}\right)}$.

Note that vector $\left(\alpha_{g}\right)_{g \in G_{0}^{+}} \in \mathbb{Q}^{\left(G_{0}^{+}\right)}$is an element of $\varphi\left(\mathcal{B}_{\text {rat }}\left(G_{0} \cup-G_{0}\right)\right)$ precisely when

$$
\sum_{g \in G_{0}^{+}} \alpha_{g} g=0 .
$$

Thus, if we let $M$ denote the $r \times\left|G_{0}^{+}\right|$matrix whose columns are the vectors $g \in G_{0}^{+} \subset \mathbb{Z}^{r} \subset \mathbb{Q}^{r}$, we see that $\varphi\left(\mathcal{B}_{\text {rat }}\left(G_{0} \cup-G_{0}\right)\right)$ is the kernel of the matrix $M$. The set $\varphi\left(\mathcal{B}_{\text {rat }}\left(G_{0}\right)\right)$ can also be described via $M$; It is the subset consisting of all vectors $\left(\alpha_{g}\right)_{g \in G_{0}^{+}} \in \operatorname{ker}(M)$ that satisfy the following sign restrictions:

$$
\alpha_{g} \geq 0 \text { unless }-g \in G_{0}
$$

(Recall that we always assume $-G_{0}^{-} \subset G_{0}^{+}$, and thus $\alpha_{g}>0$ is allowed for every $g \in G_{0}^{+}$.) It is well known that the kernel of a matrix $M$ is the orthogonal space for the row space of the same matrix $M$, and that the row and column space of $M$ have the same dimension, which in this case is equal to the dimension of the $\mathbb{Q}$-vector space spanned by the vectors from $G_{0} \subset \mathbb{Z}^{r} \subset \mathbb{Q}^{r}$. This latter number is simply $\mathrm{r}\left(\left\langle G_{0}\right\rangle\right)$ and thus we conclude that

$$
\begin{equation*}
\left|G_{0}^{+}\right|=\mathrm{r}\left(\left\langle G_{0}\right\rangle\right)+\operatorname{dim}_{\mathbb{Q}}\left(\left\langle\varphi\left(\mathcal{B}\left(G_{0}\right)\right)\right\rangle\right) . \tag{12}
\end{equation*}
$$

The next theorem (essentially due to Rockafellar [49] in a matroid formulation) shows that elementary zero-sum sequences can be useful for decomposing a zero-sum sequence via rational product decomposition. Indeed, an arbitrary zero-sum sequence always has a product decomposition into a bounded number of rational powers of elementary atoms. It also shows that an upper bound for $\mathrm{D}\left(G_{0}\right)$ can be found using an upper bound for $\mathrm{D}^{\mathrm{elm}}\left(G_{0}\right)$. It is important to note that, even if $S \in \mathcal{B}\left(G_{0}\right)$ is an atom, $S$ may still have a nontrivial product decomposition into rational powers of elementary atoms if it is not itself elementary.

Theorem 3.7. Let $r \geq 1$, let $G_{0} \subset \mathbb{Z}^{r}$ be a nonempty subset, and let $S \in \mathcal{B}_{\mathrm{rat}}\left(G_{0}\right)$ be a rational zero-sum sequence. Then there exist a nonnegative integer $\ell \geq 0$, elementary atoms $U_{1}, \ldots, U_{\ell} \in \mathcal{A}\left(G_{0}\right)$, positive rational numbers $\alpha_{0}, \alpha_{1}, \ldots, \alpha_{\ell} \in \mathbb{Q}_{>0}$, and a zero-sum sequence $R \in \operatorname{ker}(\varphi)$ (possibly trivial) having a factorization into zero-sum subsequences each of length 1 or 2 . Moreover,

$$
\begin{aligned}
& S=R^{\alpha_{0}} U_{1}^{\alpha_{1}} \cdot \ldots \cdot U_{\ell}^{\alpha_{\ell}}, \quad \operatorname{supp}^{+}\left(U_{j}\right) \nsubseteq \operatorname{supp}^{+}\left(S\left(\prod_{i=1}^{j} U_{i}^{\alpha_{i}}\right)^{-1}\right) \text { for all } i \in[1, \ell], \quad \text { and } \\
& \ell \leq \min \left\{\frac{1}{2}\left|\operatorname{supp}^{+}(S)\right|,\left|G_{0}^{+}\right|-\mathrm{r}\left(\left\langle G_{0}\right\rangle\right)\right\} .
\end{aligned}
$$

In particular,

$$
\mathrm{D}\left(G_{0}\right) \leq \sup \left\{2, \hat{\ell} \mathrm{D}^{\mathrm{elm}}\left(G_{0}\right)\right\} \leq \sup \left\{2, \min \left\{\eta,\left|G_{0}^{+}\right|-\mathrm{r}\left(\left\langle G_{0}\right\rangle\right)\right\} \mathrm{D}^{\mathrm{elm}}\left(G_{0}\right)\right\} \leq \sup \left\{2,\left|G_{0} \backslash\{0\}\right| \mathrm{D}^{\mathrm{elm}}\left(G_{0}\right)\right\}
$$

where $\eta=\sup \left\{|\operatorname{supp}(U)|: U \in \mathcal{A}\left(G_{0}\right)\right\}$ and $\hat{\ell}$ denotes the supremum over all $\ell$ as $S$ ranges over $\mathcal{A}\left(G_{0}\right)$.
Proof. We first construct the rational product decomposition for $S \in \mathcal{B}_{\text {rat }}\left(G_{0}\right)$ and then deduce the upper bound for $\mathrm{D}\left(G_{0}\right)$. To this end, in view of the observations made prior to Lemma 3.1, we may without loss of generality assume that $S$ is nontrivial and that

$$
\begin{equation*}
\left\{g \in G_{0}^{+}: g,-g \in \operatorname{supp}(S)\right\}=\emptyset \tag{13}
\end{equation*}
$$

It suffices to prove the theorem for such rational sequences.
Our goal is to show that there exist elementary atoms $U_{1}, \ldots, U_{\ell} \in \mathcal{A}\left(G_{0}\right)$ and positive rational numbers $\alpha_{1}, \ldots, \alpha_{\ell} \in \mathbb{Q}>0$ such that

$$
\begin{equation*}
S=U_{1}^{\alpha_{1}} \cdot \ldots \cdot U_{\ell}^{\alpha_{\ell}} \quad \text { and } \quad \operatorname{supp}^{+}\left(U_{j}\right) \nsubseteq \operatorname{supp}^{+}\left(S\left(\prod_{i=1}^{j} U_{i}^{\alpha_{i}}\right)^{-1}\right) \text { for all } i \in[1, \ell] \tag{14}
\end{equation*}
$$

Before doing so, we explain how (14) forces the desired upper bound for $\ell$. The second condition from (14) says that each atom $U_{j}$ contains an element $g_{j} \in \operatorname{supp}(S)$ (via the first condition of 14 p ) not contained in any of the $U_{j+1}, \ldots, U_{\ell}$. But then $\left\{g_{1}, \ldots, g_{\ell}\right\} \subset \operatorname{supp}(S)$ is a subset of cardinality $\ell$, implying $\ell \leq|\operatorname{supp}(S)|=$ $\frac{1}{2}\left|\operatorname{supp}^{+}(S)\right|$. In view of 13$)$ and since $g_{j} \in \operatorname{supp}(S)$, we see that $-g_{j} \notin \operatorname{supp}(S)$ for each $j \in[1, \ell]$. As a result, since each atom $U_{j}$ contains some element $g_{j}$ not contained in any of the $U_{j+1}, \ldots, U_{\ell}$, no $g_{j} \in \operatorname{supp}^{+}\left(U_{j}\right)$ is contained in $\operatorname{supp}^{+}\left(U_{j+1} \cdot \ldots \cdot U_{\ell}\right)$. Now it is easily deduced that $\varphi\left(U_{1}\right), \ldots, \varphi\left(U_{\ell}\right) \in \varphi\left(\mathcal{B}\left(G_{0}\right)\right)$ are linearly independent over $\mathbb{Q}$. Therefore $\ell \leq \operatorname{dim}_{\mathbb{Q}}\left(\left\langle\varphi\left(\mathcal{B}\left(G_{0}\right)\right)\right\rangle\right)=\left|G_{0}^{+}\right|-\mathrm{r}\left(\left\langle G_{0}\right\rangle\right)$, with the equality following from (12). Thus the desired bound for $\ell$ follows from (14), and we now devote our attention to proving (14). For this, suppose for the sake of contradiction that $S \in \mathcal{B}_{\text {rat }}\left(G_{0}\right)$ is a counterexample to (14) with $|\operatorname{supp}(S)|$ minimal.

Since $S$ is nontrivial, (13) forces $\operatorname{supp}^{+}(S) \neq \emptyset$. Thus, by Lemma 3.6 there exists an elementary atom $U_{1} \in \mathcal{A}\left(G_{0}\right)$ with $\operatorname{supp}\left(U_{1}\right) \subset \operatorname{supp}(S)$. Let $\alpha_{1}=\min \left\{\mathrm{v}_{g}(S) / \mathrm{v}_{g}\left(U_{1}\right): g \in \operatorname{supp}\left(U_{1}\right)\right\}$. Since $\operatorname{supp}\left(U_{1}\right) \subset$ $\operatorname{supp}(S), \alpha_{1}>0$. By definition, $\mathrm{v}_{g}\left(U_{1}^{\alpha_{1}}\right) \leq \mathrm{v}_{g}(S)$ for every $g \in \operatorname{supp}\left(U_{1}\right)$, with equality holding for some $g=g_{1} \in \operatorname{supp}\left(U_{1}\right) \subset \operatorname{supp}(S)$ (attaining the minimum in the definition of $\alpha_{1}$ ). Define $S^{\prime}=S U_{1}^{-\alpha_{1}}$. Since $\mathrm{v}_{g_{1}}\left(U_{1}^{\alpha_{1}}\right)=\mathrm{v}_{g_{1}}(S)$ with $g_{1} \in \operatorname{supp}(S)$, we conclude that $\operatorname{supp}\left(S^{\prime}\right)$ is a proper subset of $\operatorname{supp}(S)$. Indeed, $g_{1} \notin \operatorname{supp}\left(S^{\prime}\right)$. Since $S$ and $U_{1}$ are each zero-sum sequences, it follows that $S^{\prime}$ is also a zero-sum sequence. Thus, in view of the minimality of $|\operatorname{supp}(S)|$, we can apply the theorem to (the possibly trivial) rational sequence $\tilde{S}^{\prime}$ to find $\tilde{S}^{\prime}=U_{2}^{\alpha_{2}} \ldots . U_{\ell}^{\alpha_{\ell}}$ for some positive rational numbers $\alpha_{i}$ and elementary atoms $U_{i}$ satisfying (14). But now $S=U_{1}^{\alpha_{1}} S^{\prime}=U_{1}^{\alpha_{1}} U_{2}^{\alpha_{2}} \ldots . U_{\ell}^{\alpha_{\ell}}$ with (14) holding for $j \in[2, \ell]$. Since $g_{1} \in \operatorname{supp}\left(U_{1}\right)$ and $g_{1} \notin \operatorname{supp}\left(S^{\prime}\right)=\operatorname{supp}\left(U_{2} \cdot \ldots \cdot U_{\ell}\right)$, we see that 14$)$ holds when $j=1$ and thus (14) is established. This completes the proof of the first part of the theorem. It remains to prove the upper bound for $\mathrm{D}\left(G_{0}\right)$.

Let $U \in \mathcal{A}\left(G_{0}\right)$ be an atom. We must show that $|U|$ is at most the bound given at the end of Theorem 3.7. If $|U| \leq 2$, this is clearly the case and so we may assume that $|U| \geq 3$. In this case we may assume that
$\operatorname{supp}(U) \cap \operatorname{supp}(-U)=\emptyset$. Let

$$
U=U_{1}^{\alpha_{1}} \cdot \ldots \cdot U_{\ell}^{\alpha_{\ell}}
$$

be the rational product decomposition of $U$ given by the first part of the theorem. In particular, each $U_{i} \in \mathcal{A}\left(G_{0}\right)$ is an elementary atom and each $\alpha_{i} \in \mathbb{Q}_{>0}$ is a positive rational number. We note that the corresponding sequence $R$ is trivial since $\operatorname{supp}(U) \cap \operatorname{supp}(-U)=\emptyset$. If $\alpha_{i}>1$ for some $i \in[1, \ell]$, then $U_{i} \mid U$ is a proper nontrivial zero-sum subsequence, contradicting that $U \in \mathcal{A}\left(G_{0}\right)$ is an atom. Thus $\alpha_{i} \leq 1$ for all $i \in[1, \ell]$. Now

$$
|U|=\sum_{i=1}^{\ell} \alpha_{i}\left|U_{i}\right| \leq \sum_{i=1}^{\ell}\left|U_{i}\right| \leq \ell \mathrm{D}^{\mathrm{elm}}\left(G_{0}\right),
$$

and, noting that $\operatorname{supp}(U)=\frac{1}{2}\left|\operatorname{supp}^{+}(U)\right|(\operatorname{since} \operatorname{supp}(U) \cap \operatorname{supp}(-U)=\emptyset)$, the desired bound for $|U|$ follows from the bound for $\ell$ given in the first part of the theorem.

We now consider which subsets $X \subset G_{0} \subset \mathbb{Z}^{r}$ can be attained as the support of an elementary zerosum sequence. A related question asks which subsets $X \subset G_{0} \cup-G_{0} \subset \mathbb{Z}^{r}$ can be attained as the signed support of an elementary zero-sum sequence. Given any $U \in \mathcal{B}_{\text {rat }}\left(G_{0}\right)$, we know that $U^{n} \in \mathcal{B}\left(G_{0}\right)$ for some $n \geq 1$. Applying Lemma 3.5, we see that $X=\operatorname{supp}^{+}(U)$ for some elementary $U \in \mathcal{B}\left(G_{0}\right)$ is equivalent to $X=\operatorname{supp}^{+}(U)$ for some elementary $U \in \mathcal{B}_{\text {rat }}\left(G_{0}\right)$ which in turn is equivalent to $X=\operatorname{supp}^{+}(U)$ for some elementary $U \in \mathcal{A}\left(G_{0}\right)$. The same is true when considering $X=\operatorname{supp}(U)$ for an elementary zero-sum sequence $U$ over $G_{0}$. Of course, if $X=\operatorname{supp}^{+}(U)$ for a zero-sum sequence $U$, then $X \subset G_{0} \cup-G_{0}$ is symmetric and so $X=Y \cup-Y$ for some $Y \subset G_{0}$ with $Y \cap-Y=\emptyset$. The following lemma classifies the possibilities for $X$.

Lemma 3.8. Let $r \geq 1$, let $G_{0} \subset \mathbb{Z}^{r}$ be a nonempty subset, and let $X \subset G_{0} \cup-G_{0}$ be a subset with $X=-X$. Then condition (a) holds if and only if conditions (b) and (c) both hold. If $G_{0}=-G_{0}$, then (a) and (b) are equivalent.
(a) $X=\operatorname{supp}^{+}(U)$ for some elementary zero-sum sequence $U \in \mathcal{B}\left(G_{0}\right)$.
(b) The elements of $X \cap G_{0}^{+}$are linearly dependent over $\mathbb{Q}$, but any proper subset of $X \cap G_{0}^{+}$is linearly independent over $\mathbb{Q}$.
(c) There exists a nontrivial zero-sum sequence $S \in \mathcal{B}\left(G_{0}\right)$ with $\emptyset \neq \operatorname{supp}^{+}(S) \subset X$.

In particular, if $U \in \mathcal{B}\left(G_{0}\right)$ is an elementary zero-sum sequence such that $\operatorname{supp}(S) \cap \operatorname{supp}(-S)=\emptyset$, then $|\operatorname{supp}(U)| \leq r+1$.
Proof. We set $Y=X \cap G_{0}^{+}$. Suppose that (a) holds. Then (c) holds and for each $g \in Y$ there exists a nonzero $\alpha_{g} \in \mathbb{Q} \backslash\{0\}$ such that

$$
\begin{equation*}
\sum_{g \in Y} \alpha_{g} g=0 \quad \text { with } \alpha_{g}>0 \text { whenever }-g \notin G_{0} . \tag{15}
\end{equation*}
$$

Since each $\alpha_{g}$ is nonzero, the elements of $Y$ are linearly dependent.
Now, if we assume that condition (b) does not hold, then there must be some nonempty, proper subset $Z \subsetneq Y$ such that the elements of $Z$ are linearly dependent. But then for each $g \in Z$ there exist $\beta_{g} \in \mathbb{Q}$ such that

$$
\begin{equation*}
\sum_{g \in Z} \beta_{g} g=0 \tag{16}
\end{equation*}
$$

with not all $\beta_{g}$ zero. Set $\beta_{g}=0$ for any $g \in Y \backslash Z$. Suppose first that there is some $g \in Y$ with $\alpha_{g}>0$ and $\beta_{g} \neq 0$. Multiplying the $\beta_{g}$ by -1 if need be, we may assume $\alpha_{g}>0$ and $\beta_{g}>0$. In this case, let $\gamma=\min \left\{\alpha_{g} / \beta_{g}: \alpha_{g}>0, \beta_{g}>0, g \in Y\right\}>0$. If $\beta_{g}=0$ whenever $\alpha_{g}>0$, we set $\gamma=\min \left\{\alpha_{g} / \beta_{g}: \alpha_{g}<\right.$ $\left.0, \beta_{g}<0, g \in Y\right\}>0$. Multiplying (16) by $\gamma$, we obtain

$$
\begin{equation*}
\sum_{g \in Z} \gamma \beta_{g} g=0 \tag{17}
\end{equation*}
$$

Moreover, by the definition of $\gamma$ we see that $\gamma \beta_{g} \leq \alpha_{g}$ whenever $\alpha_{g}>0$ and $\beta_{g}>0$. Thus, if we subtract (17) from (15), the resulting coefficient $\alpha_{g}-\gamma \beta_{g}$ will be non-negative whenever $\alpha_{g}>0$. As a result, since $\alpha_{g}>0$ whenever $-g \notin G_{0}$ (by 15 ), we see that

$$
\begin{equation*}
\sum_{g \in Z}\left(\alpha_{g}-\gamma \beta_{g}\right) g=0 \quad \text { with } \alpha_{g}-\gamma \beta_{g} \geq 0 \text { whenever }-g \notin G_{0} \tag{18}
\end{equation*}
$$

Furthermore, for an element $g_{1} \in Y$ attaining the minimum in the definition of $\gamma$, we see that the coefficient $\alpha_{g_{1}}-\gamma \beta_{g_{1}}$ of $g_{1}$ in (18) is zero, while not all coefficients in (18) are zero since each $\alpha_{g}$ is nonzero and since at least one $\beta_{g}$ is zero ( $Z$ is a proper subset of $Y$ ). Thus the $\mathbb{Q}$-linear relation given in 18 ) corresponds to a nontrivial rational zero-sum sequence $V \in \mathcal{B}_{\text {rat }}\left(G_{0}\right)$ whose support is strictly contained in the support of $U$, contradicting that $U \in \mathcal{B}\left(G_{0}\right)$ is an elementary zero-sum sequence. We can then conclude that all proper subsets of $Y$ are linearly independent, as desired.

Now suppose that (b) holds and that either $G_{0}=-G_{0}$ or (c) holds. Clearly, there cannot be any zero-sum sequence $U$ over $G_{0} \subset \mathbb{Z}^{r} \subset \mathbb{Q}^{r}$ with $\operatorname{supp}(U)$ consisting of linearly independent elements over $\mathbb{Q}$. Thus, in order to show that (a) holds, it suffices to show that there exists some $U \in \mathcal{B}_{\text {rat }}\left(G_{0}\right)$ with $\operatorname{supp}^{+}(U)=X$.

If (c) holds, then a nontrivial $U \in \mathcal{B}_{\text {rat }}\left(G_{0}\right)$ exists with $\emptyset \neq \operatorname{supp}^{+}(U) \subset X$. However, since supp ${ }^{+}(U) \cap G^{+}$ can only be linearly independent if $\operatorname{supp}^{+}(U)$ is empty, we have $\operatorname{supp}^{+}(U)=X$ as desired.

Next assume $G_{0}=-G_{0}$. We need to show that there exists some $U \in \mathcal{B}_{\text {rat }}\left(G_{0}\right)$ with $\operatorname{supp}^{+}(U)=X$. This is equivalent to showing that there exists some nonzero $\alpha_{1}, \ldots, \alpha_{\ell} \in \mathbb{Q} \backslash\{0\}$ with $\sum_{i=1}^{\ell} \alpha_{i} g_{i}=0$. As the elements of $Y$ are linearly dependent, there exist $\alpha_{1}, \ldots, \alpha_{\ell} \in \mathbb{Q}$ not all zero with $\sum_{i=1}^{\ell} \alpha_{i} g_{i}=0$. Thus, if $\alpha_{i}=0$ for some $i \in[1, \ell]$, this would give a dependence relation on the elements of $Y \backslash\left\{g_{i}\right\}$, contradicting the hypothesis that every proper subset of $Y$ is linearly independent. Consequently, $\alpha_{i} \neq 0$ for all $i \in[1, \ell]$, and (a) follows as noted earlier.

If $U \in \mathcal{B}\left(G_{0}\right)$ is an elementary zero-sum sequence such that $\operatorname{supp}(U) \cap \operatorname{supp}(-U)=\emptyset$, then the first part of the theorem implies that $\operatorname{supp}(U) \backslash\{g\} \subset \mathbb{Q}^{r}$ is a set of linearly independent vectors for any $g \in \operatorname{supp}(U)$. Since any subset of vectors of size $r+1$ must be linearly dependent in $\mathbb{Q}^{r}$, the desired bound $|\operatorname{supp}(U)| \leq r+1$ follows.

Lemma 3.9. Let $r \geq 1$ and let $G_{0} \subset \mathbb{Z}^{r}$ be a nonempty subset. Then $\mathrm{D}\left(G_{0}\right) \geq 3$ if and only if there exists an elementary atom $U \in \mathcal{A}\left(G_{0}\right)$. If this is the case, then $\mathrm{D}^{\mathrm{elm}}\left(G_{0}\right) \geq 3$ as well.
Proof. Since any elementary atom $U \in \mathcal{A}\left(G_{0}\right)$ must satisfy $|U| \geq 3$, one direction is clear. Suppose now that $\mathrm{D}\left(G_{0}\right) \geq 3$ and let $V \in \mathcal{A}\left(G_{0}\right)$ be an atom with $|V| \geq 3$. Then $\operatorname{supp}^{+}(V)$ is nonempty, in which case Lemma 3.6 completes the proof.

Recall that, given any $m \times n$ integer matrix $M$ with $m \leq n$, we can perform elementary row and column operations on $M$ (swapping rows/columns, multiplying a row/column by $\pm 1$, or adding an integer multiple of a row/column to another row/column) to obtain a diagonal integer matrix $D=\left(d_{i, j}\right)_{i, j}$ with $d_{1,1}|\cdots| d_{m, m}$ and $d_{i, i} \geq 0$ for all $i \in[1, m]$. The matrix $D$ is unique and is known as the Smith normal form of the matrix $M$ and the $d_{i, i}$ are called the elementary divisors of $M$. If $g_{1}, \ldots, g_{n} \in \mathbb{Z}^{m}$ are the columns of $M$, then $\mathbb{Z}^{m} /\left\langle g_{1}, \ldots, g_{n}\right\rangle \cong \mathbb{Z} / d_{1,1} \mathbb{Z} \oplus \cdots \oplus \mathbb{Z} / d_{m, m} \mathbb{Z}$. Thus, when $M$ has full rank (whence $d_{m, m} \neq 0$ ), we have $d_{m, m}=\exp \left(\mathbb{Z}^{m} /\left\langle g_{1}, \ldots, g_{n}\right\rangle\right)$. It is easily checked (and well-known) that for $j \in[1, m]$

$$
\operatorname{gcd}\left\{\operatorname{det}\left(M^{\prime}\right): M^{\prime} \text { is a } j \times j \text { submatrix of } M\right\}
$$

remains invariant under elementary row and column operations and is thus equal to $d_{1,1} \cdots d_{j, j}$. In particular, if $m=n$, so that $M$ is a square matrix,

$$
\begin{equation*}
d_{m, m}=\frac{|\operatorname{det}(M)|}{\operatorname{gcd}\left\{\left|\operatorname{det}\left(M^{\prime}\right)\right|: M^{\prime} \text { is a }(m-1) \times(m-1) \text { submatrix of } M\right\}} \tag{19}
\end{equation*}
$$

These results can be found in many standard textbooks dealing with linear algebra over $\mathbb{Z}$.
We now turn our attention to finding bounds for $\mathrm{D}\left(G_{0}\right)$ where $G_{0} \subset \mathbb{Z}^{r}$. Let $M$ be a $r \times\left|G_{0}^{+}\right|$matrix whose columns are the vectors $g \in G_{0}^{+} \subset \mathbb{Z}^{r}$. Using lattice theory and results from the Geometry of Numbers, Diaconis, Graham, and Sturmfels [5] showed that

$$
\mathrm{D}\left(G_{0}\right) \leq(2 r)^{r}(r+1)^{r+1} \max \left\{\left|\operatorname{det}\left(M^{\prime}\right)\right|: M^{\prime} \text { is a } r \times r \text { submatrix of } M\right\}
$$

when $G_{0}$ is finite with full rank $r\left(\left\langle G_{0}\right\rangle\right)=r$. However, when $\left|G_{0}\right|$ is not terribly large, Theorem 3.7 can be used to obtain tighter bounds. To do so, we need to be able to bound $\mathrm{D}^{\mathrm{elm}}\left(G_{0}\right)$ and to do this we consider an argument of Sturmfels [52, Chapter 4] which, when combined with additional results, allows us to give a linear algebraic description of $\mathrm{D}^{\mathrm{elm}}\left(G_{0}\right)$. In order to state the next theorem we first need the following definition. For a collection of $r+1$ integer vectors $g_{1}, \ldots, g_{r+1} \in \mathbb{Z}^{r}$, we define

$$
\Delta\left(g_{1}, \ldots, g_{r+1}\right)=\frac{\sum_{i=1}^{r+1}\left|\operatorname{det}\left(g_{1}, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1}\right)\right|}{\operatorname{gcd}\left\{\left|\operatorname{det}\left(g_{1}, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1}\right)\right|: i \in[1, r+1]\right\}} \in \mathbb{N}_{0},
$$

where $\Delta\left(g_{1}, \ldots, g_{r+1}\right)=0$ if $\mathrm{r}\left(\left\langle g_{1}, \ldots, g_{r+1}\right\rangle\right)<r$.
Theorem 3.10. Let $r \geq 1$ and let $G_{0} \subset \mathbb{Z}^{r}$ be a nonempty subset with $\mathrm{r}\left(\left\langle G_{0}\right\rangle\right)=r$ and $\mathrm{D}\left(G_{0}\right) \geq 3$. Then

$$
\begin{aligned}
\mathrm{D}^{\mathrm{elm}}\left(G_{0}\right) & =\sup \left\{\Delta\left(g_{1}, \ldots, g_{r+1}\right): g_{1}, \ldots, g_{r+1} \in G_{0} \quad \text { and } \quad \mathrm{D}\left(\left\{g_{1}, \ldots, g_{r+1}\right\}\right) \geq 3\right\} \\
& \leq \sup \left\{\Delta\left(g_{1}, \ldots, g_{r+1}\right): g_{1}, \ldots, g_{r+1} \in G_{0}\right\} .
\end{aligned}
$$

Moreover, if $G_{0}=-G_{0}$, then equality holds.
Proof. By Lemma 3.9, the hypothesis $\mathrm{D}\left(G_{0}\right) \geq 3$ is equivalent to the existence of an elementary atom $U \in \mathcal{A}\left(G_{0}\right)$. Since any such elementary atom $U$ satisfies $|U| \geq 3$, we conclude that

$$
\begin{equation*}
\mathrm{D}^{\mathrm{elm}}\left(G_{0}\right) \geq 3 \tag{20}
\end{equation*}
$$

Let $U \in \mathcal{A}\left(G_{0}\right)$ be an arbitrary elementary atom and let $X=\operatorname{supp}(U)$. Since $U$ is an elementary atom, $\operatorname{supp}(U) \cap \operatorname{supp}(-U)=\emptyset$. By Lemma 3.8, $X$ is linearly dependent over $\mathbb{Q}$, but any proper subset of $X$ is linearly independent. In particular, $|X|=x+1 \geq 2$ with

$$
1 \leq \mathrm{r}(\langle X\rangle)=x \leq \mathrm{r}\left(\left\langle G_{0}\right\rangle\right)=r
$$

Thus, if $x<r$, then we can find a subset $X^{\prime} \subset G_{0} \backslash(X \cup-X)$ such that $\left|X^{\prime}\right|=r-x$ and $r\left(\left\langle X \cup X^{\prime}\right\rangle\right)=r$. Let $Y=X \cup X^{\prime}$. Note that $\operatorname{supp}(U) \subset Y \subset G_{0}$ with $|Y|=r+1$, and that, for each $g \in G_{0}, g$ and $-g$ are not both contained in $Y$. Let $Y=\left\{g_{1}, \ldots, g_{r+1}\right\}$ where $g_{1}, \ldots, g_{x+1}$ are the elements from $X$.

Let $M$ be the $r \times(r+1)$ matrix whose columns are the vectors $g_{i} \in Y \subset G_{0}$. Then the vector of integer multiplicities $\mathbf{x}=\left(x_{i}\right)_{i \in[1, r+1]} \in \mathbb{Z}^{|Y|}=\mathbb{Z}^{r+1}$ corresponds to a zero-sum subsequence $S=\prod_{i=1}^{r+1} g_{i}^{x_{i}} \in \mathcal{B}\left(G_{0}\right)$ with $\operatorname{supp}(S) \subset Y($ with $\operatorname{supp}(S) \cap \operatorname{supp}(-S)=\emptyset)$ and $\mathrm{v}_{g_{i}}(S)=x_{i}$ precisely when $\mathbf{x}$ is in the kernel of the matrix $M$ and $x_{i} \geq 0$ for all $i \in[1, r+1]$. Also, the vector $\mathbf{x} \in \mathbb{Z}^{r+1}$ corresponds to a zero-sum subsequence $S=\prod_{i \in I^{+}} g_{i}^{x_{i}} \prod_{i \in I^{-}}\left(-g_{i}\right)^{-x_{i}} \in \mathcal{B}\left(G_{0} \cup-G_{0}\right)$ with $\operatorname{supp}^{+}(S) \subset Y \cup-Y$ precisely when $\mathbf{x}$ is in the kernel of the matrix $M, I^{+} \subset[1, r+1]$ denotes the subset of indices $i \in[1, r+1]$ with $x_{i}>0$, and $I^{-} \subset[1, r+1]$ denotes the subset of indices $i \in[1, r+1]$ with $x_{i}<0$. In the latter case, we have $S \in \mathcal{B}\left(G_{0}\right)$ precisely when $-Y^{-} \subset G_{0}$ for $Y^{-}=\left\{y_{i} \in Y: i \in I^{-}\right\}$.

Consider the vector $\mathbf{x}=\frac{1}{\delta}\left(x_{i}\right)_{i \in[1, r+1]} \in \mathbb{Z}^{r+1}$ given by

$$
x_{i}=(-1)^{i} \operatorname{det}\left(g_{1}, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1}\right) \quad \text { for } i \in[1, r+1] \text {, }
$$

where

$$
\delta=\operatorname{gcd}\left\{\operatorname{det}\left(g_{1}, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1}\right): i \in[1, r+1]\right\} .
$$

Since $r(\langle Y\rangle)=r$, the $x_{i}$ cannot all be zero, and thus $\delta>0$ is also nonzero. Since each $g_{i} \in \mathbb{Z}^{r}$ is an integer-valued vector, it is clear from the the above definition that $\mathbf{x} \in \mathbb{Z}^{r+1}$ with

$$
\begin{equation*}
\operatorname{gcd}\left\{\frac{x_{i}}{\delta}: i \in[1, r+1]\right\}=1 \tag{21}
\end{equation*}
$$

Moreover, since $X=\left\{g_{1}, \ldots, g_{x+1}\right\}$ is linearly dependent, $x_{i}=0$ for all $i \in[x+2, r+1]$.
We now show that $\mathbf{x}$ is in the kernel of the matrix $M$ whose columns are the vectors $g_{i} \in \mathbb{Z}^{r}$. Let $g_{i}=\left(g_{i, j}\right)_{j \in[1, r]}$ with $g_{i, j}$ the $j$-th entry of the column vector $g_{i}$. With $j \in[1, r]$ arbitrary, the $j$-th entry of $M \mathbf{x} \in \mathbb{Z}^{r}$ is

$$
\begin{equation*}
\frac{1}{\delta} \sum_{i=1}^{r+1} g_{i, j} x_{i}=\frac{1}{\delta} \sum_{i=1}^{r+1}(-1)^{i} g_{i, j} \operatorname{det}\left(g_{1}, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1}\right) \tag{22}
\end{equation*}
$$

However, by the cofactor expansion formula for the determinant of a matrix, the right hand side of 22 is equal (up to sign) to the product of $\frac{1}{\delta}$ and the determinant of the $(r+1) \times(r+1)$ matrix $M^{\prime}$ formed from the $r \times(r+1)$ matrix $M$ by repeating the $j$-th row $\left(g_{i, j}\right)_{i \in[1, r]}$ of $M$ and then computing the cofactor expansion about this duplicate row. As $M^{\prime}$ has two duplicate rows, its determinant is zero and thus 22 is zero if $j \in[1, r]$. Since the $j$-th entry of $M \mathbf{x} \in \mathbb{Z}^{r}$ is equal to 22 , this shows that every coordinate of $M \mathbf{x}$ is zero. Hence $M \mathbf{x}=0$ and $\mathbf{x}$ is in the kernel of $M$ as claimed.

Since $\mathbf{x}=\frac{1}{\delta}\left(x_{i}\right)_{i \in[1, r+1]} \in \mathbb{Z}^{r+1}$ is an integer-valued vector in the kernel of $M$, it follows (as was noted earlier in the proof) that $\mathbf{x}$ corresponds to a nontrivial (not all $x_{i}$ are zero) zero-sum sequence $S \in \mathcal{B}\left(G_{0} \cup-G_{0}\right)$ with $\operatorname{supp}^{+}(S) \subset X \cup-X\left(x_{i}=0\right.$ for $i \in[x+2, r+1]$ and $\left.X=\left\{g_{1}, \ldots, g_{x+1}\right\}\right)$. Moreover, $\operatorname{supp}(S) \cap \operatorname{supp}(-S)=\emptyset$. Consequently, as $X \cup-X=\operatorname{supp}^{+}(U)$ with $U \in \mathcal{A}\left(G_{0}\right)$ elementary, $\operatorname{supp}^{+}(S)=X \cup-X$. From (21) we see that $S \neq T^{\ell}$ for any $T \in \mathcal{B}\left(G_{0} \cup-G_{0}\right)$ and $\ell \geq 2$. Since any proper subset of $\operatorname{supp}(U)$ is linearly independent, it follows that there is no nontrivial zero-sum sequence $V \in \mathcal{B}\left(G_{0} \cup-G_{0}\right)$ with $\emptyset \neq \operatorname{supp}^{+}(V) \subsetneq \operatorname{supp}^{+}(U)$. Thus $U$ is an elementary atom not just over $G_{0}$, but also over $G_{0} \cup-G_{0}$. From Lemma 3.5, $U$ must be the unique (up to sign) elementary atom over $G_{0} \cup-G_{0}$ with signed support $X \cup-X$, and all other elementary zero-sum sequences $T$ over $G_{0} \cup-G_{0}$ with $\operatorname{supp}^{+}(T)=X \cup-X$ (for which $\left.\operatorname{supp}(T) \cap \operatorname{supp}(-T)=\emptyset\right)$ must be a power of either $U$ or $-U$. Applying this conclusion to $T=S$, we find that either $S=U^{\ell}$ or $-S=U^{\ell}$ for some $\ell \geq 1$. By swapping the sign of each $x_{i}$ in the definition of $\mathbf{x}=\left(x_{i}\right)_{i \in[1, r]}$ (thus replacing $\mathbf{x}$ by $-\mathbf{x}$ ) if need be, we may without loss of generality assume the former: $S=U^{\ell}$ and, in particular, $S \in \mathcal{B}\left(G_{0}\right)$. Since $S \neq T^{\ell}$ for any $T \in \mathcal{B}\left(G_{0} \cup-G_{0}\right)$ and $\ell \geq 2$ as observed above, it follows that $\ell=1$ and $S=U$.

Now

$$
|U|=|S|=\frac{1}{\delta} \sum_{i=1}^{r}\left|x_{i}\right|=\frac{1}{\delta} \sum_{i=1}^{r+1}\left|\operatorname{det}\left(g_{1}, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1}\right)\right|=\Delta\left(g_{1}, \ldots, g_{r+1}\right) .
$$

Since $U \in \mathcal{A}\left(G_{0}\right)$ was an arbitrary elementary atom with $\operatorname{supp}(U) \subset\left\{g_{1}, \ldots, g_{r+1}\right\}$, we have

$$
\begin{align*}
\mathrm{D}^{\mathrm{elm}}\left(G_{0}\right) & \leq \sup \left\{\Delta\left(g_{1}, \ldots, g_{r+1}\right): g_{1}, \ldots, g_{r+1} \in G_{0} \quad \text { and } \quad \mathrm{D}\left(\left\{g_{1}, \ldots, g_{r+1}\right\}\right) \geq 3\right\}  \tag{23}\\
& \leq \sup \left\{\Delta\left(g_{1}, \ldots, g_{r+1}\right): g_{1}, \ldots, g_{r+1} \in G_{0}\right\} . \tag{24}
\end{align*}
$$

Let $g_{1}, \ldots, g_{r+1} \in G_{0}$ be vectors with $\mathrm{D}\left(\left\{g_{1}, \ldots, g_{r+1}\right\}\right) \geq 3$ and such that $\Delta\left(g_{1}, \ldots, g_{r+1}\right)>0$. Such vectors exist by (20) and 23). The condition $\Delta\left(g_{1}, \ldots, g_{r+1}\right)>0$ implies that $\mathrm{r}\left(\left\langle g_{1}, \ldots, g_{r+1}\right\rangle\right)=r$. Since $\mathrm{D}\left(\left\{g_{1}, \ldots, g_{r+1}\right\}\right) \geq 3$ and by Lemma 3.9. there exists an elementary atom $U \in \mathcal{A}\left(\left\{g_{1}, \ldots, g_{r+1}\right\}\right)$, such that, without loss of generality, $\operatorname{supp}(U)=\left\{g_{1}, \ldots, g_{x+1}\right\}$, where $x \leq r$. Now, repeating the above arguments using this particular elementary atom $U$, we find that $D^{\text {elm }}\left(G_{0}\right) \geq|U|=\Delta\left(g_{1}, \ldots, g_{r+1}\right)$. Taking the supremum over $\Delta\left(g_{1}, \ldots, g_{r+1}\right)$ for all choices of $g_{1}, \ldots, g_{r+1} \in G_{0}$ with $\mathrm{D}\left(\left\{g_{1}, \ldots, g_{r+1}\right\}\right) \geq 3$ and $\Delta\left(g_{1}, \ldots, g_{r+1}\right)>0$, we see that equality holds in 23 .

Next suppose that $-G_{0}=G_{0}$. To complete the proof we need to show that equality holds in (24). We may assume $\mathrm{D}^{\mathrm{elm}}\left(G_{0}\right)<\infty$ as the claim is trivially true otherwise. Let $g_{1}, \ldots, g_{r+1} \in G_{0}$ be vectors that obtain the maximum in 23). If $\mathrm{r}\left(\left\langle g_{1}, \ldots, g_{r+1}\right\rangle\right)<r$, then $\Delta\left(g_{1}, \ldots, g_{r+1}\right)=0$ and, by the lower bound $\mathrm{D}^{\text {elm }}\left(G_{0}\right) \geq 3$ from (20), the $g_{1}, \ldots, g_{r+1} \in G_{0}$ cannot maximize (24). Thus we may assume that $\mathrm{r}\left(\left\langle g_{1}, \ldots, g_{r+1}\right\rangle\right)=r$. Note
that replacing any $g_{i}$ with $-g_{i}$ does not alter the value of $\Delta\left(g_{1}, \ldots, g_{r+1}\right)$ and that the hypothesis $G_{0}=-G_{0}$ ensures that $-g_{i} \in G_{0}$. Thus, to show equality in $(24)$, it suffices to show that $\mathrm{D}\left(\left\{\epsilon_{1} g_{1}, \ldots, \epsilon_{r+1} g_{r+1}\right\}\right) \geq 3$ for some choice of $\epsilon_{i} \in\{1,-1\}$ which, by Lemma 3.9, is equivalent to the existence of an elementary atom $U \in \mathcal{A}\left(\left\{\epsilon_{1} g_{1}, \ldots, \epsilon_{r+1} g_{r+1}\right\}\right)$ for some choice of $\epsilon_{i} \in\{1,-1\}$ which in turn is equivalent to the existence of an elementary atom $U \in \mathcal{A}\left(G_{0}\right)$ with $\operatorname{supp}^{+}(U) \subset\left\{g_{1}, \ldots, g_{r+1}\right\} \cup-\left\{g_{1}, \ldots, g_{r+1}\right\}$. To show the later we will use Lemma 3.8.

Let $Y \subset\left\{g_{1}, \ldots, g_{r+1}\right\}$ be those $g_{i} \in\left\{g_{1}, \ldots, g_{r+1}\right\}$ such that $\operatorname{det}\left(g_{1}, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1}\right) \neq 0$. It follows that $\left\{g_{1}, \ldots, g_{r+1}\right\} \backslash\left\{g_{i}\right\}$ is linearly independent for any $g_{i} \in Y$. Thus, any subset of $\left\{g_{1}, \ldots, g_{r+1}\right\} \backslash$ $\left\{g_{i}\right\}$, including $Y \backslash\left\{g_{i}\right\}$, must also be linearly independent. This shows that all proper subsets of $Y$ are linearly independent.

Suppose the elements of $Y$ are linearly independent. Then clearly $|Y| \leq r$ and thus $Y$ must be a proper subset of $\left\{g_{1}, \ldots, g_{r+1}\right\}$. However, since $\mathrm{r}\left(\left\langle g_{1}, \ldots, g_{r+1}\right\rangle\right)=r$, it follows that we can complete $Y$ to some full rank set $Y^{\prime} \subset\left\{g_{1}, \ldots, g_{r+1}\right\}$ with $Y \subset Y^{\prime}$. Then $\operatorname{det}\left(g_{1}, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1}\right) \neq 0$ for the unique $g_{i} \in\left\{g_{1}, \ldots, g_{r+1}\right\} \backslash Y^{\prime}$. By the definition of $Y$, this forces $g_{i} \in Y$, contradicting that $g_{i} \notin Y^{\prime}$ with $Y \subset Y^{\prime}$. Thus we conclude that the elements of $Y$ are linearly dependent.

In view of the conclusions of the previous two paragraphs, along with the hypothesis $G_{0}=-G_{0}$, we can now apply Lemma 3.8 and conclude that there exists an elementary zero-sum sequence $S \in \mathcal{B}\left(G_{0}\right)$ with $\operatorname{supp}^{+}(S)=Y \subset\left\{g_{1}, \ldots, g_{r+1}\right\} \cup-\left\{g_{1}, \ldots, g_{r+1}\right\}$. By Lemma 3.5. this ensures that there exists an elementary atom $U \in \mathcal{A}\left(G_{0}\right)$ with $\operatorname{supp}^{+}(U)=Y \subset\left\{g_{1}, \ldots, g_{r+1}\right\} \cup-\left\{g_{1}, \ldots, g_{r+1}\right\}$ which completes the proof.

Corollary 3.11. Let $r \in \mathbb{N}$, let $\left(\mathrm{e}_{1}, \ldots, \mathrm{e}_{r+1}\right)$ denote the standard basis of $\mathbb{Z}^{r+1}$, and let $G_{0} \subset G=\left\langle\mathrm{e}_{1}, \ldots, \mathrm{e}_{r}\right\rangle$ be a nonempty subset with $\mathrm{r}\left(\left\langle G_{0}\right\rangle\right)=r$ and $\mathrm{D}\left(G_{0}\right) \geq 3$. Furthermore, let $\widetilde{G}_{0}, \widetilde{G}_{1} \subset \mathbb{Z}^{r+1}$ be the subsets given by

$$
\begin{aligned}
\widetilde{G}_{0} & =\left\{g+\mathrm{e}_{r+1}: g \in G_{0}\right\} \cup G_{0} \quad \text { and } \\
\widetilde{G}_{1} & =\left\{g+\mathrm{e}_{r+1}: g \in G_{0}\right\} \cup\left\{g-\mathrm{e}_{r+1}: g \in G_{0}\right\},
\end{aligned}
$$

and let $\mathcal{M}\left(\widetilde{G}_{i}\right)$ be all those nonsingular $(r+1) \times(r+1)$ matrices with columns $\tilde{g}_{1}, \ldots, \tilde{g}_{r+1} \in \widetilde{G}_{i} \subset \mathbb{Z}^{r+1}$, for $i \in[0,1]$. Also, let $d_{r+1}(\tilde{M})$ denote the largest elementary divisor of the matrix $\tilde{M}$.

1. $\mathrm{D}^{\operatorname{elm}}\left(G_{0}\right) \leq 2 \sup \left\{d_{r+1}(\tilde{M}): \tilde{M} \in \mathcal{M}\left(\widetilde{G}_{0}\right)\right\} \leq 2 \sup \left\{|\operatorname{det}(\tilde{M})|: \tilde{M} \in \mathcal{M}\left(\widetilde{G}_{0}\right)\right\}$ and
2. $\mathrm{D}^{\mathrm{elm}}\left(G_{0}\right) \leq \sup \left\{d_{r+1}(\tilde{M}): \tilde{M} \in \mathcal{M}\left(\widetilde{G}_{1}\right)\right\} \leq \sup \left\{|\operatorname{det}(\tilde{M})|: \tilde{M} \in \mathcal{M}\left(\widetilde{G}_{1}\right)\right\}$.

Proof. By Lemma 3.9 and Theorem 3.10 we know that

$$
3 \leq \mathrm{D}^{\mathrm{elm}}\left(G_{0}\right) \leq \sup \left\{\Delta\left(g_{1}, \ldots, g_{r+1}\right): g_{1}, \ldots, g_{r+1} \in G_{0}\right\}
$$

It follows that there exist $g_{1}, \ldots, g_{r+1} \in G_{0} \subset \mathbb{Z}^{r}$ with $\Delta\left(g_{1}, \ldots, g_{r+1}\right)>0$ and that the supremum on the right hand side is necessarily obtained for such a choice of $g_{1}, \ldots, g_{r+1}$. Let $g_{1}, \ldots, g_{r+1} \in G_{0}$ be such that $\Delta\left(g_{1}, \ldots, g_{r+1}\right)>0$. Then $\mathrm{r}\left(\left\langle g_{1}, \ldots, g_{r+1}\right\rangle\right)=r$.

For each $i \in[1, r+1]$, let $\tilde{g}_{i}=g_{i} \pm \mathrm{e}_{r+1}$, where an appropriate choice for the sign of $\mathrm{e}_{r+1}$ will be determined shortly, and let $\tilde{M} \in \mathcal{M}\left(\widetilde{G}_{1}\right)$ be the $(r+1) \times(r+1)$ integer matrix with columns $\tilde{g}_{1}, \ldots, \tilde{g}_{r+1} \in \tilde{G}_{1} \subset \mathbb{Z}^{r+1}$. Let

$$
\begin{aligned}
& \delta=\operatorname{gcd}\left\{\left|\operatorname{det}\left(g_{1}, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1}\right)\right|: i \in[1, r+1]\right\} \quad \text { and } \\
& \tilde{\delta}=\operatorname{gcd}\left\{\left|\operatorname{det}\left(M^{\prime}\right)\right|: M^{\prime} \text { is a } r \times r \text { sub-matrix of } \tilde{M}\right\} .
\end{aligned}
$$

Then $\tilde{\delta} \mid \delta$. By using the cofactor expansion formula for the determinant and expanding along the final row of $\tilde{M}$, we see that by choosing the signs in each $\tilde{g}_{i}$ appropriately, $\frac{1}{\delta} \operatorname{det}(\tilde{M})=\Delta\left(g_{1}, \ldots, g_{r+1}\right)$. Now, since $\tilde{\delta} \mid \delta$ and by 19), we have

$$
\Delta\left(g_{1}, \ldots, g_{r+1}\right)=\frac{1}{\delta} \operatorname{det}(\tilde{M}) \leq \frac{1}{\tilde{\delta}} \operatorname{det}(\tilde{M})=d_{r+1}(\tilde{M})
$$

Since the choice of $g_{1}, \ldots, g_{r+1} \in G_{0}$ was arbitrary, subject to the restriction $\Delta\left(g_{1}, \ldots, g_{r+1}\right)>0$, the bound from Theorem 3.10 establishes 2.

For each $i \in[1, r+1]$, redefine $\tilde{g}_{i}$ as either $\tilde{g}_{i}=g_{i}+\mathrm{e}_{r+1}$ or $\tilde{g}_{i}=g_{i}$, where the choice of coefficient for $\mathrm{e}_{r+1}$ will be determined shortly, and let $\tilde{M} \in \mathcal{M}\left(\widetilde{G}_{0}\right)$ be the $(r+1) \times(r+1)$ integer matrix with columns $\tilde{g}_{1}, \ldots, \tilde{g}_{r+1} \in \tilde{G}_{0} \subset \mathbb{Z}^{r+1}$. Let

$$
\begin{aligned}
& \delta=\operatorname{gcd}\left\{\left|\operatorname{det}\left(g_{1}, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1}\right)\right|: i \in[1, r+1]\right\} \quad \text { and } \\
& \tilde{\delta}=\operatorname{gcd}\left\{\left|\operatorname{det}\left(M^{\prime}\right)\right|: M^{\prime} \text { is a } r \times r \text { sub-matrix of } \tilde{M}\right\} .
\end{aligned}
$$

Then $\tilde{\delta} \mid \delta$. By again using the cofactor expansion formula for the determinant and expanding along the final row of $\tilde{M}$, we see that by choosing the coefficients for $\mathrm{e}_{r+1}$ in each $\tilde{g}_{i}$ appropriately, we can achieve

$$
\operatorname{det}(\tilde{M})=\sum_{i \in I}\left|\operatorname{det}\left(g_{1}, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1}\right)\right|
$$

for some subset $I \subset[1, r]$. By using the exact opposite choices for the $\tilde{g}_{i}$, we can instead achieve

$$
\operatorname{det}(\tilde{M})=-\sum_{i \in[1, r+1] \backslash I}\left|\operatorname{det}\left(g_{1}, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1}\right)\right|
$$

Between these two options, take the one where $|\operatorname{det}(\tilde{M})|$ is larger. Then

$$
2|\operatorname{det}(\tilde{M})| \geq \sum_{i=1}^{r+1}\left|\operatorname{det}\left(g_{1}, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1}\right)\right|
$$

in which case (19) and $\tilde{\delta} \mid \delta$ give

$$
\Delta\left(g_{1}, \ldots, g_{r+1}\right) \leq \frac{2}{\delta}|\operatorname{det}(\tilde{M})| \leq \frac{2}{\tilde{\delta}}|\operatorname{det}(\tilde{M})|=2 d_{r+1}(\tilde{M})
$$

establishing 1.
The special case $G_{0}=G_{r}^{+} \cup-G_{r}^{+}$. Our goal for the remainder of this section is to apply the machinery above to the case when

$$
G_{0}=G_{r}^{+} \cup-G_{r}^{+}
$$

where $G_{r}^{+}$denotes the nonzero vertices of the $r$-dimensional hypercube as defined in Section 2 . Specifically, we wish to obtain upper and lower bounds for $\mathrm{D}\left(G_{r}^{+} \cup-G_{r}^{+}\right)$which we will apply to the study of invariants of monoids of modules in Section 4.

If $r=1$, then clearly $\mathrm{D}\left(G_{0}\right)=2$. Thus we suppose for the remainder of this section that $r>1$ in which case $\mathrm{D}\left(G_{0}\right) \geq 3$. Note that $\Delta\left(g_{1}, \ldots, g_{r+1}\right)$ is unaffected if any $g_{i}$ is replaced by $-g_{i}$. Thus we need only consider matrices with columns from $G_{r}^{+}$when applying Theorem 3.10 and Corollary 3.11. By Corollary 3.11 we see that $\mathrm{D}^{\mathrm{elm}}\left(G_{0}\right)$ is bounded from above by twice the maximal determinant of a $(r+1) \times(r+1)(0,1)$ matrix. It is well-known that this value is $\frac{1}{2^{r+1}}$ times the maximal value of a $(r+2) \times(r+2)(1,-1)$-matrix and this value, in turn, is equal to the maximal value of a $(r+2) \times(r+2)(-1,0,1)$-matrix. Moreover, this maximal is bounded by Hadamard's bound and, together with Theorem 3.7, we obtain the following upper bound

$$
\mathrm{D}\left(G_{0}\right) \leq\left(\left|G_{0}^{+}\right|-r\right) \mathrm{D}^{\mathrm{elm}}\left(G_{0}\right) \leq \frac{\left(2^{r}-r-1\right)}{2^{r}}(r+2)^{\frac{r+2}{2}}<(r+2)^{\frac{r+2}{2}}
$$

We now construct a lower bound for $\mathrm{D}\left(G_{0}\right)$. Suppose we have an $(r+1) \times(r+1)(0,1)$-matrix $M$ whose determinant is a large prime $p$. Then we must have $d_{r+1}=p$ and $d_{r}=1$ where the $d_{1}|\cdots| d_{r+1}$ are the elementary divisors of $M$. Therefore, by $\sqrt{19}$, we see that the greatest common divisor of the determinants of $r \times r$ sub-matrices of $M$ must be 1 . On the other hand, if we delete any row of $M$, we obtain, by applying the cofactor expansion formula to $\operatorname{det}(M)$ and expanding along the deleted, row $r+1$ vectors $g_{1}, \ldots, g_{r+1} \in G_{0}$
where $\operatorname{gcd}\left\{\left|\operatorname{det}\left(g_{1}, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1}\right)\right|: \quad i \in[1, r+1]\right\} \operatorname{divides} \operatorname{det}(M)=p$. As $p$ is prime, we can conclude that $\operatorname{gcd}\left\{\operatorname{det}\left(g_{1}, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1}\right): i \in[1, r+1]\right\}$ is either 1 or $p$. Moreover, since the greatest common divisor of the determinants of $r \times r$ sub-matrices is 1 , we see that we can achieve 1 rather than $p$ by choosing an appropriate row of $M$ to delete. Then, in this case, $\Delta\left(g_{1}, \ldots, g_{r+1}\right)=d_{r+1}=|\operatorname{det}(M)|=p$ and we obtain the lower bound

$$
\max \{|\operatorname{det}(M)|: M \text { is a }(r+1) \times(r+1) \quad(0,1) \text {-matrix with }|\operatorname{det}(M)| \operatorname{prime}\} \leq \mathrm{D}^{\mathrm{elm}}\left(G_{0}\right) \leq \mathrm{D}\left(G_{0}\right) .
$$

In general, the possible values of $|\operatorname{det}(M)|$ for an arbitrary $(0,1)$-matrix are not known, though this question is of great interest to many researchers (see [46] for known results on the spectrum of the determinant). However, computational evidence obtained for small values of $r$ has led many to observe that there is (at least for $r$ small) a constant $C \cong \frac{1}{2}$ and a large consecutive interval $\left[0, C 2^{-r-1}(r+2)^{(r+2) / 2}\right]$ of obtainable values for $|\operatorname{det}(M)|$, leading some to conjecture that the set of values of $\operatorname{det}(M)$ is dense in a interval whose length is a fraction of the maximal possible value (see [45). For any interval of obtainable values $[0, n]$, Bertrand's postulate ensures that a prime of size at least $\frac{1}{2} n$ will occur in that interval. Thus, if the intuitions gathered from examining small values of $r$ hold true for larger values in a very strong sense, we would expect a lower bound for $\mathrm{D}\left(G_{0}\right)$ of the form

$$
\left(\frac{r+2}{C}\right)^{(r+2) / 2} \lesssim \mathrm{D}^{\mathrm{elm}}\left(G_{0}\right) \leq \mathrm{D}\left(G_{0}\right)
$$

where $C \geq 1$ is some constant. We note that such a lower bound would very nearly match the upper bound in order of magnitude.

Constructively, the best lower bounds we have been able to achieve involve the Fibonacci numbers. We now present a construction giving a lower bound on $\mathrm{D}\left(G_{0}\right)$ regardless of concerns about the possible values of $\operatorname{det}(M)$. We note that for small values of $r$, examples of $(0,1)$-matrices with large prime determinant are known and can thus be used to improve upon this bound. For $r \in \mathbb{N}_{0}$, we denote by $\boldsymbol{F}_{r}$ the $r$ th Fibonacci number. That is, $\mathrm{F}_{0}=0, \mathrm{~F}_{1}=1$, and $\mathrm{F}_{r}=\mathrm{F}_{r-1}+\mathrm{F}_{r-2}$ for all $r \geq 2$.

Proposition 3.12. Let $r \in \mathbb{N}$, let $\left(\mathrm{e}_{1}, \ldots, \mathrm{e}_{r}\right)$ denote the standard basis of $\mathbb{Z}^{r}$, and set $H_{r}=\left\langle\mathrm{e}_{1}+\cdots+\mathrm{e}_{r}\right\rangle$. Then there exists a sequence $S_{r} \in \mathcal{F}\left(G_{r}^{+}\right)$with

$$
\sigma\left(S_{r}\right) \in H_{r} \quad \text { and } \quad \Sigma_{\leq\left|S_{r}\right|-1}\left(S_{r}\right) \cap H_{r}=\emptyset
$$

such that

$$
\begin{array}{ll}
\left|\operatorname{supp}\left(S_{r}\right)\right|=r & \text { with }
\end{array} \quad \operatorname{supp}\left(S_{r}\right) \text { spanning } \mathbb{Q}^{r}, ~ 子, ~\left(S_{r}\right)=\mathrm{F}_{r} \mathrm{e}_{1}+\cdots+\mathrm{F}_{r} \mathrm{e}_{r} .
$$

Proof. The sequence $S_{1}=\mathrm{e}_{1}$ is easily seen to satisfy the conditions of the theorem. We proceed recursively to define $S_{r}$ for $r \geq 2$, assuming that

$$
S_{r-1}=g_{1} \cdot \ldots \cdot g_{\mathrm{F}_{r}} \in \mathcal{F}\left(G_{r-1}^{+}\right)
$$

has already been constructed so as to satisfy the conclusions of the theorem.
Let

$$
S_{r}^{\prime}=S_{r-1} \mathrm{e}_{r}^{\mathrm{F}_{r-1}} \in \mathcal{F}\left(G_{r}^{+}\right)
$$

and let

$$
S_{r}=\varphi\left(S_{r}^{\prime}\right) \in \mathcal{F}\left(\mathbb{Z}^{r}\right),
$$

where $\varphi: \mathbb{Z}^{r} \rightarrow \mathbb{Z}^{r}$ is the map defined by $x \mapsto-\left(x-\left(\mathrm{e}_{1}+\cdots+\mathrm{e}_{r}\right)\right)$ and extends to a affine linear isomorphism of $\mathbb{Q}^{r}$. The map $\varphi$ acts on an element $x \in G_{r}^{+}$simply by exchanging each 0 for a 1 and each 1 for a 0 . Therefore $\varphi(x) \in G_{r}^{+}$for each $x \in G_{r}^{+} \backslash\left\{\mathrm{e}_{1}+\cdots+\mathrm{e}_{r}\right\}$. As a result, since $\mathrm{e}_{1}+\cdots+\mathrm{e}_{r} \notin \operatorname{supp}\left(S_{r}^{\prime}\right)\left(\operatorname{as} \operatorname{supp}\left(S_{r-1}\right) \subset G_{r-1}^{+}\right.$ and $r \geq 2$ ), we see that

$$
S_{r} \in \mathcal{F}\left(G_{r}^{+}\right) \quad \text { with } \quad\left|S_{r}\right|=\left|S_{r}^{\prime}\right| \quad \text { and } \quad\left|\operatorname{supp}\left(S_{r}\right)\right|=\left|\operatorname{supp}\left(S_{r}^{\prime}\right)\right| .
$$

Observe that

$$
\left|S_{r}\right|=\left|S_{r}^{\prime}\right|=\left|S_{r-1}\right|+\mathrm{F}_{r-1}=\mathrm{F}_{r}+\mathrm{F}_{r-1}=\mathrm{F}_{r+1}
$$

and that

$$
\begin{aligned}
\sigma\left(S_{r}\right) & =\sigma\left(\varphi\left(S_{r}^{\prime}\right)\right)=\left|S_{r}\right| \mathrm{e}_{1}+\cdots+\left|S_{r}\right| \mathrm{e}_{r}-\sigma\left(S_{r}^{\prime}\right) \\
& =\mathrm{F}_{r+1} \mathrm{e}_{1}+\cdots+\mathrm{F}_{r+1} \mathrm{e}_{r}-\left(\mathrm{F}_{r-1} \mathrm{e}_{1}+\cdots+\mathrm{F}_{r-1} \mathrm{e}_{r}\right) \\
& =\mathrm{F}_{r} \mathrm{e}_{1}+\cdots+\mathrm{F}_{r} \mathrm{e}_{r} \in H_{r} .
\end{aligned}
$$

Moreover, since $\varphi: \mathbb{Q}^{r} \rightarrow \mathbb{Q}^{r}$ is a affine linear isomorphism, we have that

$$
\left|\operatorname{supp}\left(S_{r}\right)\right|=\left|\operatorname{supp}\left(\varphi\left(S_{r}^{\prime}\right)\right)\right|=\left|\operatorname{supp}\left(S_{r}^{\prime}\right)\right|=\left|\operatorname{supp}\left(S_{r-1}\right)\right|+1=r .
$$

Since $\varphi: \mathbb{Q}^{r} \rightarrow \mathbb{Q}^{r}$ is a affine linear isomorphism, to show that $\operatorname{supp}\left(S_{r}\right)=\operatorname{supp}\left(\varphi\left(S_{r}^{\prime}\right)\right)$ spans $\mathbb{Q}^{r}$ it suffices to show that $\operatorname{supp}\left(S_{r}^{\prime}\right)$ spans $\mathbb{Q}^{r}$. But this is clear since $\operatorname{supp}\left(S_{r}^{\prime}\right)=\operatorname{supp}\left(S_{r-1}\right) \cup\left\{\mathrm{e}_{r}\right\}$ with $\operatorname{supp}\left(S_{r-1}\right)$ spanning $\mathbb{Q} e_{1}+\cdots+\mathbb{Q} e_{r-1}$ by hypothesis. It remains to show that $\Sigma_{\leq\left|S_{r}\right|-1}\left(S_{r}\right) \cap H_{r}=\emptyset$. Since $\sigma(\varphi(T)) \in H_{r}$ if and only if $\sigma(T) \in H_{r}$ for any sequence $T \in \mathcal{F}\left(\mathbb{Z}^{r}\right)$, we see that in order to show $\Sigma_{\leq\left|S_{r}\right|-1}\left(S_{r}\right) \cap H_{r}=\emptyset$, it suffices to show that $\Sigma_{\leq\left|S_{r}^{\prime}\right|-1}\left(S_{r}^{\prime}\right) \cap H_{r}=\emptyset$.

Suppose that $T \mid S_{r}$ is a nontrivial subsequence with $\sigma(T) \in H_{r}$, that is, the coordinates of each entry of $T$ are equal. Then the coordinates of the first $r-1$ entries are equal. However, since $S_{r}^{\prime}=S_{r-1} \mathbf{e}_{r}^{\mathbf{F}_{r-1}}$ and by the hypothesis that $\Sigma_{\leq\left|S_{r-1}\right|-1}\left(S_{r-1}\right) \cap H_{r-1}=\emptyset$, this is only possible if either $T \mid \mathrm{e}_{r}^{\mathbf{F}_{r-1}}$ or $S_{r-1} \mid T$. In the former case, since $r \geq \overline{2}$, it is clear that $\sigma(T) \notin H$. In the latter case, since $\sigma\left(S_{r-1}\right)=\mathrm{F}_{r-1} \mathrm{e}_{1}+\cdots+\mathrm{F}_{r-1} \mathrm{e}_{r-1}$ and $S_{r}^{\prime}=S_{r-1} \mathrm{e}_{r}^{\mathbf{F}_{r-1}}, \sigma(T) \in H$ only if $T=S_{r}^{\prime}$. Thus $\Sigma_{\leq\left|S_{r}^{\prime}\right|-1}\left(S_{r}^{\prime}\right) \cap H_{r}=\emptyset$ follows, completing the proof.

Theorem 3.13. Let $r \in \mathbb{N}_{\geq 2}$, let $\left(\mathrm{e}_{1}, \ldots, \mathrm{e}_{r}\right)$ denote the standard basis of $\mathbb{Z}^{r}$, and let $G_{0}=G_{r}^{+} \cup-G_{r}^{+}$. Then

$$
\mathrm{F}_{r+2} \leq \mathrm{D}\left(G_{0}\right) \leq \frac{\left(2^{r}-r-1\right)}{2^{r}}(r+2)^{\frac{r+2}{2}} \leq(r+2)^{\frac{r+2}{2}} .
$$

Proof. The upper bounds follow from Corollary 3.11 (see the discussion following the corollary). It remains to show $\mathrm{F}_{r+2} \leq \mathrm{D}\left(G_{0}\right)$. Let $S_{0} \in \mathcal{F}\left(G_{r}^{+}\right)$be a sequence satisfying the conclusions of Proposition 3.12 and define

$$
U=S_{r}\left(-\mathrm{e}_{1}-\cdots-\mathrm{e}_{r}\right)^{\mathbf{F}_{r}} \in \mathcal{F}\left(G_{0}\right) .
$$

Then $|U|=\left|S_{r}\right|+\mathrm{F}_{r}=\mathrm{F}_{r+1}+\mathrm{F}_{r}=\mathrm{F}_{r+2}$ and $\sigma(U)=\sigma\left(S_{r}\right)-\left(\mathrm{F}_{r} \mathrm{e}_{1}+\cdots+\mathrm{F}_{r} \mathrm{e}_{r}\right)=0$. Moreover, by the definition of $U$ it is clear that if $T \mid U$ is a zero-sum subsequence with $T=T^{+} T^{-}$where $\operatorname{supp}\left(T^{+}\right) \subset G_{r}^{+}$ and $\operatorname{supp}\left(T^{-}\right) \subset-G_{r}^{+}$, then we must have $\sigma\left(T^{+}\right) \in H$ with $T^{+} \mid S_{r}$. However, since $\sigma\left(S_{r}\right) \in H_{r}$ and $\Sigma_{\leq\left|S_{r}\right|-1}\left(S_{r}\right) \cap H_{r}=\emptyset$, this is only possible if either $T^{+}$is trivial or $T^{+}=S_{r}$. If $T^{+}$is trivial, then $T=\left(-\mathrm{e}_{1}-\cdots-\mathrm{e}_{r}\right)^{|T|}$ which is a zero-sum sequence only if $T$ is trivial. If $T^{+}=S_{r}$, then $\sigma\left(T^{+}\right)=\sigma\left(S_{r}\right)=$ $\mathrm{F}_{r} \mathrm{e}_{1}+\cdots+\mathrm{F}_{r} \mathrm{e}_{r}$ and it then follows from the definition of $U$ that the only way $T$ can be a zero-sum is if $T=U$. Therefore $U$ is a zero-sum sequence of length $\mathrm{F}_{r+2}$ having no proper nontrivial zero-sum subsequences. Thus $\mathrm{D}\left(G_{0}\right) \geq \mathrm{F}_{r+2}$ as desired.

Remark 3.14. Let $r \in \mathbb{N},\left(\mathrm{e}_{1}, \ldots, \mathrm{e}_{r}\right)$, and $G_{0}$ be as in Theorem 3.13 but with $r>1$. From Theorem 3.7 we know that $\mathrm{D}^{\mathrm{elm}}\left(G_{0}\right) \leq \mathrm{D}\left(G_{0}\right) \leq \frac{1}{2}\left|G_{0}\right| \mathrm{D}^{\mathrm{elm}}\left(G_{0}\right)$. But the expected value of $\mathrm{D}^{\mathrm{elm}}\left(G_{0}\right)$ is much larger than $\left|G_{0}\right|$, meaning that $\mathrm{D}\left(G_{0}\right) \approx \mathrm{D}^{\mathrm{elm}}\left(G_{0}\right)$. Indeed, by a computer based search we have verified that $\mathrm{D}\left(G_{0}\right)=\mathrm{D}^{\mathrm{elm}}\left(G_{0}\right)$ for $r \in\{2,3\}$. However, whether $\mathrm{D}\left(G_{0}\right)=\mathrm{D}^{\mathrm{elm}}\left(G_{0}\right)$ remains true for $r \geq 4$ is not known. For the module-theoretic relevance of this question see [2, Lemma 6.9 and Corollary 6.10].

## 4. Monoids of modules over commutative Noetherian local Rings

In this section we study direct-sum decompositions of finitely generated modules over commutative Noetherian local rings. With $R$ a commutative Noetherian local ring and $\mathcal{C}$ a class of finitely generated modules over $R$, it is well-known that the monoid of modules $\mathcal{V}(\mathcal{C})$ is Krull. Our first lemma summarizes some basic information about $\mathcal{C}$.

Lemma 4.1. Let $(R, \mathfrak{m})$ be a commutative Noetherian local ring with maximal ideal $\mathfrak{m}$ and let $(\widehat{R}, \widehat{\mathfrak{m}})$ denote its $\mathfrak{m}$-adic completion. Let $\mathcal{C}$ denote the class of all finitely generated $R$-modules and let $\mathcal{C}^{\prime}$ denote a subclass of $\mathcal{C}$ such that $\mathcal{V}\left(\mathcal{C}^{\prime}\right) \subset \mathcal{V}(\mathcal{C})$ is a divisor-closed submonoid.

1. The embedding $\mathcal{V}(\mathcal{C}) \hookrightarrow \mathcal{V}(\widehat{\mathcal{C}})$, defined by $[M] \mapsto[M \otimes \widehat{R}]$, is a divisor homomorphism into the free abelian monoid $\mathcal{V}(\widehat{\mathcal{C}})$, where $\widehat{\mathcal{C}}$ denotes the class of finitely generated $\widehat{R}$-modules. In particular, $\mathcal{V}(\mathcal{C})$ is a Krull monoid.
2. $\mathcal{V}\left(\mathcal{C}^{\prime}\right)$ is a Krull monoid whose class group is an epimorphic image of a subgroup of the class group of $\mathcal{V}(\mathcal{C})$. If $\mathcal{V}\left(\mathcal{C}^{\prime}\right)$ is tame, then each of the arithmetical finiteness results of Proposition 2.1 hold.
3. Suppose, in addition, that $R$ is one-dimensional and reduced (no non-zero nilpotent elements). Let $G$ denote the class group of $\mathcal{V}\left(\mathcal{C}^{\prime}\right)$ and let $G_{P} \subset G$ denote the set of classes containing prime divisors.
(a) The class group $G$ of $\mathcal{V}\left(\mathcal{C}^{\prime}\right)$ is a finitely generated abelian group.
(b) $\mathcal{V}\left(\mathcal{C}^{\prime}\right)$ is tame if and only if $\mathrm{D}\left(G_{P}\right)<\infty$ if and only if $G_{P}$ is finite.

Proof. The monoid $\mathcal{V}(\widehat{\mathcal{C}})$ is free abelian by the Theorem of Krull-Remak-Schmidt-Azumaya. Also, the embedding defined by $[M] \mapsto[M \otimes \hat{R}]$ is a divisor homomorphism by [53] or [40, Corollary 1.15], and hence $\mathcal{V}(\mathcal{C})$ is a Krull monoid, proving 1.

Let $H=\mathcal{V}(\mathcal{C})$ and suppose that the inclusion $H \hookrightarrow F=\mathcal{F}(P)$ a divisor theory. Then $\mathcal{C}(H)=\mathrm{q}(F) / \mathrm{q}(H)$. If $H^{\prime} \subset H$ is divisor-closed, the inclusion $H^{\prime} \hookrightarrow F^{\prime}=\mathcal{F}\left(P^{\prime}\right)$, where $P^{\prime}=\{p \in P: p$ divides some $a \in H^{\prime}$ in $\left.F\right\}$, is a divisor homomorphism. Note that $H^{\prime}=F^{\prime} \cap \mathrm{q}\left(H^{\prime}\right)$ and that $\mathrm{q}\left(H^{\prime}\right)=\mathrm{q}\left(F^{\prime}\right) \cap \mathrm{q}(H)$. Therefore the homomorphism $\mathrm{q}\left(F^{\prime}\right) \rightarrow \mathrm{q}(F) / \mathrm{q}(H)$, defined by $a \mapsto a \mathrm{q}(H)$ for each $a \in \mathrm{q}\left(F^{\prime}\right)$, has kernel $\mathrm{q}\left(F^{\prime}\right) \cap \mathrm{q}(H)=\mathrm{q}\left(H^{\prime}\right)$. Consequently, there exists a monomorphism $\mathrm{q}\left(F^{\prime}\right) / \mathrm{q}\left(H^{\prime}\right) \rightarrow \mathrm{q}(F) / \mathrm{q}(H)=\mathcal{C}(H)$. Finally, [28, Theorem 2.4.7] implies that the class group of $H^{\prime}$ is an epimorphic image of a subgroup of $\mathrm{q}\left(F^{\prime}\right) / \mathrm{q}\left(H^{\prime}\right)$, and hence of a subgroup of $\mathcal{C}(H)$, proving 2 .

We first consider the first statement of 3 . The class group of $\mathcal{V}(\mathcal{C})$ is free abelian of finite rank by 38, Theorem 6.3]. Therefore, by statement 2 , the class group $G$ of $\mathcal{V}\left(\mathcal{C}^{\prime}\right)$ is finitely generated. We now consider the second statement of 3 . Since $G$ is finitely generated, [30, Theorem 4.2] implies that $\mathcal{V}\left(\mathcal{C}^{\prime}\right)$ is tame if and only if $\mathrm{D}\left(G_{P}\right)<\infty$. Again, since $G$ is finitely generated, [28, Theorem 3.4.2] implies that $\mathrm{D}\left(G_{P}\right)<\infty$ if and only if $G_{P}$ is finite.

Remark 4.2. We now make two brief remarks on certain hypothesis in Statement 3 of Lemma 4.1.

1. The assumption that $R$ is reduced can be slightly weakened (see [38, Section 6]). However, the assumption that $R$ is one-dimensional is essential for guaranteeing that the class group of $\mathcal{V}\left(\mathcal{C}^{\prime}\right)$ is finitely generated (see [40, Lemma 2.16]). Indeed, there is a two-dimensional complete Noetherian local Krull domain $S$ whose class group is not finitely generated. By a result of Heitman, there is a factorial two-dimensional Noetherian local domain $R$ whose completion is isomorphic to $S$. Then, if $\mathcal{C}$ is the class of finitely generated torsion-free $R$-modules, the class group of $\mathcal{V}(\mathcal{C})$ is isomorphic to the class group of $S$ (see [2, Section 5]), and hence not finitely generated.
2. Both characterizations in $3(\mathrm{~b})$ strongly depend on the fact that the class group of $\mathcal{V}\left(\mathcal{C}^{\prime}\right)$ is finitely generated and thus the hypothesis that $R$ is one-dimensional in $3(\mathrm{a})$ is critical for the results of this section pertaining to monoids of modules.

Let $R$ be a one-dimensional reduced commutative Noetherian local ring and let $\mathcal{C}^{\prime}$ be a class of finitely generated $R$-modules such that $\mathcal{V}\left(\mathcal{C}^{\prime}\right)$ is a divisor-closed submonoid of the monoid $\mathcal{V}(\mathcal{C})$ of all finitely generated $R$-modules. Then $\mathcal{V}\left(\mathcal{C}^{\prime}\right)$ is a Krull monoid with finitely generated class group $G$ and we let $G_{P} \subset G$ denote the set of classes containing prime divisors. By Proposition 2.3, sets of lengths in $\mathcal{V}\left(\mathcal{C}^{\prime}\right)$ can be studied in the monoid $\mathcal{B}\left(G_{P}\right)$ of zero-sum sequences over $G_{P}$. Using the combinatorial results of Section 3, we study sets of lengths in such Krull monoids in Theorem 4.3 in the case where $G_{P}$ contains the set $G_{r}^{+} \cup G_{r}^{-}$where $G_{r}^{+}$ denotes the nonzero vertices of the hypercube in a finitely generated free abelian group. We then investigate finer arithmetical invariants of $\mathcal{V}\left(\mathcal{C}^{\prime}\right)$ for small values of $r$. We conclude this section with an explicit example of a monoid of modules realizing such a Krull monoid (see Corollary 4.7).

Theorem 4.3. Let $H$ be a Krull monoid whose class group $\mathcal{C}(H)$ is free abelian with basis $\left(e_{1}, \ldots, e_{r}\right)$ for some $r \geq 2$, and let $G_{P} \subset \mathcal{C}(H)$ denote the set of classes containing prime divisors. Suppose that $G_{P}$ is finite and that $G_{P}=-G_{P}$ with $G_{P} \supset G_{r}^{+}$.

1. There exists $M \in \mathbb{N}_{0}$ such that, for each $k \geq 2, \mathcal{U}_{k}(H)=L^{\prime} \cup L^{*} \cup L^{\prime \prime}$, where $L^{*}$ is an interval, $L^{\prime} \subset \min L^{*}+[-M,-1]$, and $L^{\prime \prime} \subset \max L^{*}+[1, M]$.
2. For each $k \in \mathbb{N}$,

$$
\rho_{2 k}(H)=k \mathrm{D}\left(G_{P}\right) \geq k \mathrm{~F}_{r+2} \quad \text { and } \quad k \mathrm{D}\left(G_{P}\right)+1 \leq \rho_{2 k+1}(H) \leq k \mathrm{D}\left(G_{P}\right)+\frac{\mathrm{D}\left(G_{P}\right)}{2}
$$

Upper bounds for $\mathrm{D}\left(G_{P}\right)$, and hence for each $\rho_{k}(H)$, then follow from Theorem 3.13.
Proof. By Proposition 2.3, $\mathcal{B}\left(G_{P}\right)$ and $H$ are tame and thus, by Lemma 2.2 , it suffices to prove all assertions about $H$ for the monoid $\mathcal{B}\left(G_{P}\right)$.

By Proposition 2.1.3, it suffices to show that $\min \Delta\left(G_{P}\right)=1$. If $U=e_{1} e_{2}\left(-e_{1}-e_{2}\right)$, then $U \in \mathcal{A}\left(G_{P}\right)$ with $\mathrm{L}((-U) U)=\{2,3\}$. Therefore $1 \in \Delta\left(G_{P}\right)$, proving 1 .

Since $G_{P} \supset G_{r}^{+} \cup-G_{r}^{+}$, the monoid $\mathcal{B}\left(G_{r}^{+} \cup-G_{r}^{+}\right)$is a divisor closed submonoid of $\mathcal{B}\left(G_{P}\right)$ and thus $\mathrm{D}\left(G_{P}\right) \geq \mathrm{D}\left(G_{r}^{+} \cup-G_{r}^{+}\right)$. Theorem 3.13 now implies that $\mathrm{D}\left(G_{P}\right) \geq \mathrm{D}\left(G_{r}^{+} \cup-G_{r}^{+}\right) \geq \mathrm{F}_{r+2}$. The inequalities involving $\rho_{k}\left(G_{P}\right)$ now follow easily from the definitions. Indeed, if $k \in \mathbb{N}$ and $A \in \mathcal{B}\left(G_{P}\right)$ with

$$
A=0^{\mathrm{v}_{0}(A)} U_{1} \cdot \ldots \cdot U_{k}=0^{\mathrm{v}_{0}(A)} V_{1} \cdot \ldots \cdot V_{l}
$$

where $U_{1}, \ldots, U_{k}, V_{1}, \ldots, V_{l} \in \mathcal{A}\left(G_{P}\right) \backslash\{0\}$ are minimal zero-sum sequences, then

$$
2 l \leq \sum_{i=1}^{l}\left|V_{i}\right|=|A|-\mathrm{v}_{0}(A)=\sum_{i=1}^{k}\left|U_{i}\right| \leq k \mathrm{D}\left(G_{P}\right)
$$

whence $l \leq k \mathrm{D}\left(G_{P}\right) / 2$ and thus $\rho_{k}\left(G_{P}\right) \leq k \mathrm{D}\left(G_{P}\right) / 2$. If $U=g_{1} \cdot \ldots \cdot g_{l}$ is a minimal zero-sum sequence of length $|U|=l=\mathrm{D}\left(G_{P}\right)$, then

$$
\begin{equation*}
U^{k}(-U)^{k}=\prod_{i=1}^{l}\left(\left(-g_{i}\right) g_{i}\right)^{k} \tag{25}
\end{equation*}
$$

and hence $\rho_{2 k}\left(G_{P}\right) \geq k \mathrm{D}\left(G_{P}\right)$. Multiplying each side of 25) by any fixed atom, we find that $\rho_{2 k+1}\left(G_{P}\right) \geq$ $k \mathrm{D}\left(G_{P}\right)+1$.

We now make a conjecture that claims that, at least for sufficiently large $k \in \mathbb{N}$, the first statement of Theorem 4.3 holds with $M=0$.

Conjecture 4.4. Let $H$ be a Krull monoid as in Theorem 4.3 and suppose, in addition, that $G_{P}=G_{r}^{+} \cup-G_{r}^{+}$. Then there exists $k^{*} \in \mathbb{N}$ such that for each $k \geq k^{*}, \mathcal{U}_{k}(H)$ is an interval.

We now discuss one possible strategy for proving Conjecture 4.4. If one could show that there exists $A^{*} \in \mathcal{B}\left(G_{P}\right)$ with $\mathrm{L}\left(A^{*}\right)$ an interval and $\max \mathrm{L}\left(A^{*}\right) / \min \mathrm{L}\left(A^{*}\right)=\mathrm{D}\left(G_{P}\right) / 2$, then there must exist $k^{*} \in \mathbb{N}$ such that for each $k \geq k^{*}, \mathcal{U}_{k}(H)$ is an interval ([20, Theorem 3.1]). Unfortunately, this strategy seems to require knowledge of the precise value of the Davenport constant, which is currently known only for $r \in[2,3]$. Even for $r=4$, it seems to be computational infeasible to compute the Davenport constant. However, for small values of $r$ we can provide a direct proof of Conjecture 4.4. Indeed, we are even able to show in Corollary 4.6 that the conjecture holds when $r \in[2,3]$ for $k^{*}=2$. We first provide a simple lemma.

Lemma 4.5. Let $G$ be a free abelian group of rank $r \in \mathbb{N}$, let $\left(e_{1}, \ldots, e_{r}\right)$ denote a basis for $G$, and set $G_{0}=G_{r}^{+} \cup-G_{r}^{+}$. For every $U \in \mathcal{A}\left(G_{0}\right)$ with $|U| \geq 3$ and any $g \in G_{0}$, we have $\mathrm{v}_{g}(U)<|U| / 2$.

Proof. Let $U=g^{k} g_{1} \cdot \ldots \cdot g_{l} \in \mathcal{A}\left(G_{0}\right)$ where $k=\mathrm{v}_{g}(U)$ and $|U|=k+l$. By symmetry we may suppose that $g \in G_{r}^{+}$. Clearly, we have $k \leq l=|U|-k$ and hence $k \leq|U| / 2$. Assume for the sake of contradiction that $2 k=|U|$. Then $g_{1}, \ldots, g_{l} \in-G_{r}^{+}$. For each $h \in G_{0}$ with $h=\sum_{i=1}^{r} a_{i} e_{i}$ for some $a_{1}, \ldots, a_{r} \in \mathbb{Z}$, we denote by $S(h)$ the set $S(h)=\left\{i \in[1, r]: a_{i} \neq 0\right\}$. We now obtain that

$$
S\left(g_{1}\right) \cup \cdots \cup S\left(g_{l}\right) \subset S(g) \subset \bigcap_{i=1}^{l} S\left(g_{i}\right),
$$

and hence $g_{1}=\cdots=g_{l}=-g$. Unless $U=g(-g)$, this contradicts the fact that $U \in \mathcal{A}\left(G_{0}\right)$.

Corollary 4.6. Let $H$ be a Krull monoid whose class group $G$ is free abelian with basis ( $e_{1}, \ldots, e_{r}$ ) for some $r \in[2,3]$. Let $G_{P} \subset G$ denote the set of classes containing prime divisors and suppose that $G_{P}=G_{r}^{+} \cup-G_{r}^{+}$.

1. Suppose $r=2$. Then $\mathrm{c}(H)=\mathrm{t}(H)=\mathrm{D}\left(G_{P}\right)=3$ and, for each $k \geq 2$, the set $\mathcal{U}_{k}(H)$ is an interval. Moreover, for all $k \geq 2$ and $j \in[0,1], \rho_{2 k+j}(H)=3 k+j$.
2. Suppose $r=3$. Then $\mathrm{c}(H)=\mathrm{D}\left(G_{P}\right)=5$ and, for each $k \geq 2$, the set $\mathcal{U}_{k}(H)$ is an interval. Moreover, for all $k \geq 2$ and $j \in[0,1], \rho_{2 k+j}(H)=5 k+j$.

Proof. When $r=2$, the statement $\mathrm{c}(H)=\mathrm{t}(H)=\mathrm{D}\left(G_{P}\right)=3$ is a simple observation. Indeed, this setting is a special case of [2, Theorem 6.4]. The remaining assertions follow from Theorem 4.3.

We now assume $r=3$. A lengthy technical proof or a computer search show that $\mathrm{D}\left(G_{P}\right)=5$ and that $\left\{V \in \mathcal{A}\left(G_{P}\right):|V|=5\right\}=\left\{V_{i},-V_{i}: i \in[1,4]\right\}$ where

$$
\begin{aligned}
& V_{1}=\left(e_{1}+e_{2}\right)\left(e_{1}+e_{3}\right)\left(e_{2}+e_{3}\right)\left(-e_{1}-e_{2}-e_{3}\right)^{2}, \\
& V_{2}=\left(e_{1}+e_{2}\right)\left(e_{1}+e_{3}\right)\left(-e_{2}-e_{3}\right)\left(-e_{1}\right)^{2}, \\
& V_{3}=\left(-e_{1}-e_{3}\right)\left(e_{1}+e_{2}\right)\left(e_{2}+e_{3}\right)\left(-e_{2}\right)^{2}, \quad \text { and } \\
& V_{4}=\left(-e_{1}-e_{2}\right)\left(e_{1}+e_{3}\right)\left(e_{2}+e_{3}\right)\left(-e_{3}\right)^{2} .
\end{aligned}
$$

Lemma 4.5 then implies that each minimal zero-sum sequence of length four is squarefree. Clearly, $\mathrm{L}\left(\left(-V_{1}\right) V_{1}\right)=$ $\{2,5\}$. Proposition 2.3 now implies that $5 \leq 2+\max \Delta(H) \leq \mathrm{c}(H) \leq \mathrm{D}\left(G_{P}\right)=5$, whence $\mathrm{c}(H)=5$. Again, by Proposition 2.3, it suffices to prove the remaining assertions for $\mathcal{B}\left(G_{P}\right)$.

Theorem 4.3 implies that $\rho_{2 k}\left(G_{P}\right)=k \mathrm{D}\left(G_{P}\right)=5 k$ and that $5 k+1 \leq \rho_{2 k+1}\left(G_{P}\right) \leq 5 k+2$ for all $k \in \mathbb{N}$. In order to prove that $\mathcal{U}_{k}(H)$ is an interval for each $k \in \mathbb{N}$, it suffices to show that $\mathcal{U}_{k}\left(G_{P}\right) \cap \mathbb{N}_{\geq k}$ is an interval for each $k \in \mathbb{N}$. Indeed, this follows from a simple symmetry argument (see [20, Lemma 3.5]). We proceed conclude the proof by proving the following three claims A1, A2, and A3 which clearly imply the assertion.

A1. For each $k \in \mathbb{N}, \mathcal{U}_{2 k}\left(G_{P}\right) \cap \mathbb{N}_{\geq 2 k}=[2 k, 5 k]$.
A2. For each $k \in \mathbb{N}, \rho_{2 k+1}\left(G_{P}\right)=5 k+1$.
A3. For each $k \in \mathbb{N}, \mathcal{U}_{2 k+1}\left(G_{P}\right) \cap \mathbb{N}_{\geq 2 k+1}=[2 k+1,5 k+1]$.

Proof of A1. For any $k \in \mathbb{N}$, the inclusion $\mathcal{U}_{2 k}\left(G_{P}\right) \cap \mathbb{N}_{\geq 2 k} \subset\left[2 k, \rho_{2 k}\left(G_{P}\right)\right]=[2 k, 5 k]$ is clear. To prove the reverse inclusion we proceed by induction on $k$. For every $j \in[3,5]$ there exists $U \in \mathcal{A}\left(G_{P}\right)$ with $|U|=j$, and it follows that $\{2, j\} \subset \mathrm{L}((-U) U)$. Therefore $\mathcal{U}_{2}\left(G_{P}\right)=[2,5]$ and, together with the induction hypothesis, we see that for $k \geq 2$,

$$
[2 k, 5 k]=[2,5]+[2 k-2,5 k-5]=\mathcal{U}_{2}\left(G_{P}\right)+\mathcal{U}_{2 k-2}\left(G_{P}\right) \subset \mathcal{U}_{2 k}\left(G_{P}\right) .
$$

Proof of A2. Assume for the sake of contradiction that there exists $k^{\prime} \in \mathbb{N}$ such that $\rho_{2 k^{\prime}+1}\left(G_{P}\right)=5 k^{\prime}+2$, and let $k \in \mathbb{N}$ denote the smallest integer with this property. Let $B \in \mathcal{B}\left(G_{P}\right)$ be such that

$$
B=A_{1} \cdot \ldots \cdot A_{2 k+1}=B_{1} \cdot \ldots \cdot B_{5 k+2}, \quad \text { where } \quad A_{1}, \ldots, A_{2 k+1}, B_{1}, \ldots, B_{5 k+2} \in \mathcal{A}\left(G_{P}\right) .
$$

Then

$$
10 k+4 \leq \sum_{\nu=1}^{5 k+2}\left|B_{\nu}\right|=|B|=\sum_{\nu=1}^{2 k+1}\left|A_{\nu}\right| \leq 5(2 k+1)
$$

and hence, after renumbering if necessary, $\left|A_{1}\right|=\cdots=\left|A_{2 k}\right|=5, A_{2 k+1} \in[4,5],\left|B_{1}\right|=\cdots=\left|B_{5 k+1}\right|=2$, and $\left|B_{5 k+2}\right| \in[2,3]$. Suppose there are $i, j \in[1,2 k+1]$ with $i<j$ such that $A_{i}=-A_{j}$. Without loss of generality, $i=1$ and $j=2$ in which case $B^{\prime}=A_{3} \cdot \ldots \cdot A_{2 k+1}$ satisfies $2 k-1,5(k-1)+2 \in \mathrm{~L}\left(B^{\prime}\right)$, contradicting the minimality of $k$.

Thus there are distinct $U_{1}, \ldots, U_{4} \in\left\{V_{i},-V_{i}: i \in[1,4]\right\}$ such that $\left\{A_{i}: i \in[1,2 k]\right\} \subset\left\{U_{1}, \ldots, U_{4}\right\}$ with $U_{i} \neq-U_{j}$ for all $i, j \in[1,4]$. Suppose that $U_{1}=g_{1}^{2} g_{2} g_{3} g_{4}$ such that $\left|\left\{i \in[1,2 k]: A_{i}=U_{1}\right\}\right| \geq\lceil k / 2\rceil$. Then $\mathrm{v}_{g_{1}}(B) \geq \max \{k, 2\}$. By inspection of all $W \in \mathcal{A}\left(G_{P}\right)$ with $|W|=5$, it follows that $\mathrm{v}_{-g_{1}}\left(U_{1} \cdot \ldots \cdot U_{4}\right)=0=$ $\mathrm{v}_{-g_{1}}\left(A_{1} \cdot \ldots \cdot A_{2 k}\right)$, and hence $\mathrm{v}_{-g_{1}}(B)=\mathrm{v}_{-g_{1}}\left(A_{2 k+1}\right)$. Since $\left|B_{1}\right|=\cdots=\left|B_{5 k+1}\right|=2,\left|B_{5 k+2}\right| \in[2,3]$, and $\mathrm{v}_{g_{1}}\left(B_{5 k+2}\right) \leq 1$, it follows that

$$
\mathrm{v}_{-g_{1}}\left(A_{2 k+1}\right)=\mathrm{v}_{-g_{1}}(B) \geq \mathrm{v}_{g_{1}}(B)-1 \geq \max \{k-1,1\},
$$

and thus $\left|A_{2 k+1}\right|=4$. Therefore $\left|B_{5 k+2}\right|=2, \mathrm{v}_{-g_{1}}\left(A_{2 k+1}\right)=\mathrm{v}_{-g_{1}}(B)=\mathrm{v}_{g_{1}}(B) \geq \max \{k, 2\}$. However, as we noted before, each minimal zero-sum sequences of length four is squarefree, a contradiction.

Proof of A3. Let $k \in \mathbb{N}$. By assertion A2, $\mathcal{U}_{2 k+1}\left(G_{P}\right) \cap \mathbb{N}_{\geq 2 k+1} \subset\left[2 k+1, \rho_{2 k+1}\left(G_{P}\right)\right]=[2 k+1,5 k+1]$. On the other hand,

$$
[2 k+1,5 k+1]=1+[2 k, 5 k]=\mathcal{U}_{1}\left(G_{P}\right)+\mathcal{U}_{2 k}\left(G_{P}\right) \subset \mathcal{U}_{2 k+1}\left(G_{P}\right)
$$

Module theory provides an abundance of examples of Krull monoids satisfying the assumptions of Theorem 4.3 (see [2, Section 4], or [40, Chapter 1] for various realization results). Below we provide one specific example of a Krull monoid where the set $G_{P}$ of classes containing prime divisors is precisely $G_{P}=G_{r}^{+} \cup G_{r}^{-}$where $G_{r}$ is the set of nonzero vertices of the $r$-dimensional hypercube.

Corollary 4.7. Let $(R, \mathfrak{m})$ be a one-dimensional analytically unramified commutative Noetherian local domain with unique maximal ideal $\mathfrak{m}$. Further assume that the $\mathfrak{m}$-adic completion $\widehat{R}$ of $R$ has $r+1$ minimal primes $\mathfrak{q}_{1}, \ldots, \mathfrak{q}_{r+1}$, and let $M$ be a torsion-free $R$-module whose completion $\widehat{M}=M \otimes_{R} \widehat{R}$ satisfies

$$
\widehat{M} \cong \bigoplus_{\emptyset \neq I \subset[1, r+1]} \frac{\widehat{R}}{\cap_{i \in I} \mathfrak{q}_{i}}
$$

Then $\mathcal{V}(\operatorname{add}(M))$ is a Krull monoid with class group $G \cong \mathbb{Z}^{r}$ and the set $G_{P}$ of classes containing prime divisors satisfies $G_{P}=G_{r}^{+} \cup-G_{r}^{+}$, where $G, G_{P}$, and $G_{r}^{+}$are as in Theorem 4.3. Therefore all arithmetical invariants, including all sets $\mathcal{U}_{k}(\mathcal{V}(\operatorname{add}(M)))$, are as described in Theorem 4.3 and Corollary 4.6.

Proof. The statements about $G$ and $G_{P}$ follow from [2, Example 4.21]. The arithmetical consequences then follow from Theorem 4.3 and Corollary 4.6 .

## 5. Monoids of modules over Prüfer Rings

In this section we study classes $\mathcal{C}$ of finitely presented modules over Prüfer rings and characterize the algebraic structure of the monoid $\mathcal{V}(\mathcal{C})$. Specifically, we study certain classes of projective modules over various types of Prüfer rings, and show that they are always half-factorial. We also study the catenary and tame degrees of these monoids. We first recall the definition of a Prüfer ring and related topics as well as that of finitely primary monoids. For a general reference on modules over Prüfer rings, the reader may wish to consult the monograph of Fuchs and Salce [22. For modules over Prüfer rings with zero-divisors, we refer the reader to 18. For additional information on finitely primary monoids, see [28, Sections 2.7, 3.1, and 4.3].

A Prüfer ring is a commutative ring in which every finitely generated regular ideal is invertible. A commutative ring $R$ has

- the $1 \frac{1}{2}$ generator property if, for any invertible ideal $I \subset R$ and any regular element $a \in I \backslash \operatorname{rad}(R) I$, there exists an element $b \in R$ such that $I=R a+R b$.
- small zero-divisors if for every zero-divisor $a \in R$ and any ideal $A \subset R, A+a R=R$ implies that $A=R$.

A monoid $H$ is called finitely primary if there exist $s, \alpha \in \mathbb{N}$ with the following properties:
$H$ is a submonoid of a factorial monoid $F=F^{\times} \times\left[p_{1}, \ldots, p_{s}\right]$ for $s$ pairwise non-associated prime elements $p_{1}, \ldots, p_{s}$ satisfying

$$
H \backslash H^{\times} \subset p_{1} \cdot \ldots \cdot p_{s} F \quad \text { and } \quad\left(p_{1} \cdot \ldots \cdot p_{s}\right)^{\alpha} F \subset H
$$

In this case we say that $H$ is finitely primary of rank $s$ and exponent $\alpha$.
It is easy to show that the complete integral closure of such a finitely primary monoid $H$ is $F$, and hence $H$ is a Krull monoid if and only if $H$ is factorial. Moreover, $H_{\text {red }}$ is finitely generated if and only if $s=1$ and $\left(\widehat{H}^{\times}: H^{\times}\right)<\infty$. The main (and motivating) examples of finitely primary monoids stem from ring theory. For example, if $R$ is a one-dimensional local Mori domain with nonzero conductor ( $R: \widehat{R}$ ) and $\widehat{R}$ denotes the complete integral closure of $R$, then $R^{\bullet}$ is finitely primary ( $[28$, Proposition 2.10.7]). The arithmetic of finitely primary monoids is well-studied ( [28, Sections 2.7, 3.1, and 4.3]). In particular, the sets $\mathcal{U}_{k}(H)$ are finite (for one $k \geq 2$, or equivalently for all $k \geq 2$ ) if and only if $s=1$. In our main results of this section, Theorems 5.1 and 5.3 , we apply these arithmetical results to monoids of modules of modules over Prüfer rings.

Theorem 5.1. Let $R$ be a Prüfer ring such that $R$ has the $1 \frac{1}{2}$ generator property and $R$ has small zerodivisors. Let $\mathcal{C}_{\text {proj }}$ be the class of finitely generated projective $R$-modules.

1. $\mathcal{V}\left(\mathcal{C}_{\text {proj }}\right)$ is a finitely primary monoid of rank 1 and of exponent 1 . Moreover, $\mathcal{V}\left(\mathcal{C}_{\text {proj }}\right)$ is finitely generated if and only if $\operatorname{Pic}(R)$ is finite.
2. If $\mathcal{C}_{\text {proj }}$ does not satisfy KRSA, then $\mathcal{V}\left(\mathcal{C}_{\text {proj }}\right)$ is half-factorial with $\mathrm{c}\left(\mathcal{V}\left(\mathcal{C}_{\text {proj }}\right)\right)=\mathrm{t}\left(\mathcal{V}\left(\mathcal{C}_{\text {proj }}\right)\right)=2$.

Proof. We first consider statement 1. By [18, Corollary 4], every module $P$ in $\mathcal{C}_{\text {proj }}$ is isomorphic to $R^{n-1} \oplus I$ where $n \in \mathbb{N}$ is the rank of $P, I$ is an invertible ideal, and the isomorphism class of $P$ is determined by that of $I$ and by the rank $n$. Thus the map $\varphi: \mathcal{V}\left(\mathcal{C}_{\text {proj }}\right) \rightarrow H=(\operatorname{Pic}(R) \times \mathbb{N}) \cup\{(0,0)\}$, defined by $P \mapsto([I], n)$, is an isomorphism. By definition, $H \subset \operatorname{Pic}(R) \times\left(\mathbb{N}_{0},+\right)$ is finitely primary of rank 1 and exponent 1 . Since $H$ is reduced, we obtain that $\left(\widehat{H}^{\times}: H^{\times}\right)=|\operatorname{Pic}(R)|$ and thus $H$ is finitely generated if and only if $\operatorname{Pic}(R)$ is finite.

We now compute the catenary and tame degrees based on the monoid described in 1 . Suppose that $\mathcal{C}_{\text {proj }}$ does not satisfy KRSA. Then $H$ is not factorial and hence $\mathrm{c}(H) \geq 2$. By [28, Theorem 3.1.5], every finitely primary monoid of rank 1 and exponent $\alpha$ satisfies $\mathrm{c}(H) \leq \mathrm{t}(H) \leq 3 \alpha-1$, and hence $\mathrm{c}(H)=\mathrm{t}(H)=2$. Now Proposition 2.1 implies that $\Delta(H)=\emptyset$, that is, $H$ is half-factorial.

We now restrict our attention to direct-sum decompositions of modules over Prüfer domains. Specifically, we consider the class of all finitely presented modules, including torsion modules. Before proceeding, we recall the following characterization of Prüfer domains. An integral domain $R$ is Prüfer if and only if the following two equivalent conditions are satisfied:
(P1) The torsion submodule of any finitely generated $R$-module $M$ is isomorphic to a direct summand of $M$ ([22, Chap.V, Cor. 2.9]).
(P2) Every finitely generated $R$-module is projective if and only if it is torsion-free ([22, Chap.V, Theorem 2.7]).

In Theorem 5.3 we can provide a more precise arithmetical description of $\mathcal{V}(\mathcal{C})$, where $\mathcal{C}$ is the class of finitely presented modules over a Prüfer domain, if we further assume that the domain is $h$-local. We now recall this class of integral domains. Let $R$ be a domain and, for an ideal $I \subset R$, let $\Omega(I)$ denote the set of maximal ideals of $R$ containing $I$. Note, that $|\Omega(I)|=1$ implies that $R / I$ is local and hence indecomposable as an $R$-module. Also recall that a domain $R$ has finite character if each nonzero element of $R$ is contained in at most finitely many maximal ideals of $R$. Now, we say that an integral domain $R$ is $h$-local if the following equivalent conditions are satisfied ([44, Theorem 2.1]):
(H1) $R$ has finite character and each nonzero prime ideal of R is contained in a unique maximal ideal.
(H2) For each nonzero ideal $I$ of $R, R / I$ has a decomposition $\oplus_{\nu=1}^{m} R / I_{\nu}$ with $\left|\Omega\left(I_{1}\right)\right|=\cdots=\left|\Omega\left(I_{m}\right)\right|=1$.
(H3) Each torsion torsion $R$-module $M$ is canonically isomorphic to $\oplus_{\mathfrak{p} \in \max (R)} M_{\mathfrak{p}}$.
For additional information on $h$-local domains we refer the reader to the survey article by Olberding 44] and to the monograph by Fontana, Houston, and Lucas [19]. The next proposition gathers together the module-theoretic results necessary for the arithmetical results we present in Theorem 5.3. We would like to thank Bruce Olberding for the short proof of Proposition 5.2.1.

Proposition 5.2. Let $R$ be a commutative ring.

1. Let $m \in \mathbb{N}$ and let $I, I_{1}, \ldots, I_{m}$ be ideals of $R$ such that $R / I \cong R / I_{1} \oplus \cdots \oplus R / I_{m}$. If $I$ is finitely generated and projective as an $R$-module, then $I_{1}, \ldots, I_{m}$ are also finitely generated and projective as $R$-modules.
2. Let $m, n \in \mathbb{N}$ and let $I_{1}, \ldots, I_{m}, J_{1}, \ldots, J_{n}$ be ideals of $R$. If $I_{m} \subset \cdots \subset I_{1}, J_{n} \subset \cdots \subset J_{1}$, and

$$
R / I_{1} \oplus \cdots \oplus R / I_{m} \cong R / J_{1} \oplus \cdots \oplus R / J_{n}
$$

then $m=n$ and $I_{\nu}=J_{\nu}$ for each $\nu \in[1, m]$.
3. Let $R$ be an h-local Prüfer domain and let $M$ be a finitely presented $R$-module. Then, as an $R$-module, $M$ decomposes as

$$
M \cong R / I_{1} \oplus \cdots \oplus R / I_{m} \oplus I_{m+1} \oplus \cdots \oplus I_{n}
$$

where $n \in \mathbb{N}_{0}, m \in[0, n], I_{m} \subset \cdots \subset I_{1}$ are proper invertible ideals of $R$, and $I_{m+1}, \ldots, I_{n}$ are invertible ideals of $R$.

Proof. Clearly, it is sufficient to prove that $I_{1}$ is finitely generated and projective. We set $J=I_{2} \cap \cdots \cap I_{m}$ and, by [22, Lemma 1.1, Chap. V], obtain that $I=I_{1} \cap J$ and $R=I_{1}+J$. Therefore there is a short exact sequence

$$
0 \rightarrow I \rightarrow I_{1} \oplus J \rightarrow R \rightarrow 0
$$

of $R$-modules where the second map is the embedding and where the third map is given by $(x, y) \mapsto x-y$ for all $x \in I_{1}$ and all $y \in J$. This sequence splits, and thus $I_{1} \oplus J \cong R \oplus I$. Now, since $I$ is finitely generated and projective, so is $I_{1}$. This proves 1 .

For the proofs of statements 2 and 3, see Proposition 2.10 and Theorem 4.12 in [22, Chap. V].
We now state our main results about finitely presented modules over Prüfer domains.

Theorem 5.3. Let $R$ be a Prüfer domain, $\mathcal{C}$ the class of all finitely presented $R$-modules, $\mathcal{C}_{\text {tor }}$ the class of finitely presented torsion modules, and $\mathcal{C}_{\text {proj }}$ the class of finitely generated projective modules.

1. $\mathcal{V}(\mathcal{C})=\mathcal{V}\left(\mathcal{C}_{\text {proj }}\right) \times \mathcal{V}\left(\mathcal{C}_{\text {tor }}\right)$.
2. If $R$ has the $1 \frac{1}{2}$ generator property, then $\mathcal{V}\left(\mathcal{C}_{\text {proj }}\right)$ is finitely primary of rank 1 and exponent 1 .
3. Assume, in addition, that $R$ is h-local. Then $\mathcal{V}\left(\mathcal{C}_{\text {tor }}\right)$ is free abelian and $\mathcal{V}(\mathcal{C})$ is half-factorial. Furthermore, $\mathcal{V}(\mathcal{C})$ is either factorial or $\mathrm{c}(\mathcal{V}(\mathcal{C}))=\mathrm{t}(\mathcal{V}(\mathcal{C}))=2$.
Proof. The proof of 1 follows immediately from (P1) and (P2), and statement 2 follows immediately from Theorem 5.1

We now prove statement 3 , and we begin by showing that $\mathcal{V}\left(\mathcal{C}_{\text {tor }}\right)$ is free abelian. Since $R$ has finite character, $R$ has the $1 \frac{1}{2}$ generator property. Let $M$ be a finitely presented non-zero torsion $R$-module. We now argue that the module $M$ has a decomposition as a direct sum of indecomposable finitely presented $R$-modules, and that such a decomposition is unique up to isomorphism. By Proposition 5.2,3, M has a decomposition

$$
M \cong R / I_{1} \oplus \cdots \oplus R / I_{m}
$$

with $m \in \mathbb{N}$ and invertible ideals $I_{m} \subset \cdots \subset I_{1} \subsetneq R$. Property (H2) then implies that for each $\nu \in[1, m]$, $R / I_{\nu}$ has a direct-sum decomposition into indecomposables modules, each of the form $R / J$ where $J \subset R$ and $|\Omega(J)|=1$. Since each $I_{\nu}$ is invertible, each $I_{\nu}$ is finitely generated and projective and, by Proposition 5.2. 1, the same is true for all ideals $J$ with $R / J$ occurring in the direct sum decomposition of $R / I_{\nu}$. Thus, after replacing the $R / I_{\nu}$ with direct-sums of finitely generated indecomposable $R$-modules of the form $R / J$ and then renaming, we may suppose that each $R / I_{\nu}$ is an indecomposable $R$-module, that each $I_{\nu}$ is an invertible ideal, and that $\left|\Omega\left(I_{\nu}\right)\right|=1$ for each $\nu \in[1, m]$.

Let $M \cong C_{1} \oplus \cdots \oplus C_{n}$ be any direct-sum decomposition of $M$ into indecomposable finitely presented $R$-modules. Then Proposition 5.2 and (H2) imply that, for each $\nu \in[1, n], C_{\nu} \cong R / J_{\nu}$ for some invertible ideal $J_{\nu} \subset R$ with $\left|\Omega\left(J_{\nu}\right)\right|=1$. Therefore $M \cong R / J_{1} \oplus \cdots \oplus R / J_{n}$. Let $\mathfrak{p}$ be a maximal ideal of $R$. Then

$$
M_{\mathfrak{p}} \cong\left(R / I_{1}\right)_{\mathfrak{p}} \oplus \cdots \oplus\left(R / I_{m}\right)_{\mathfrak{p}} \cong\left(R / J_{1}\right)_{\mathfrak{p}} \oplus \cdots \oplus\left(R / J_{n}\right)_{\mathfrak{p}}
$$

and, by (H3), it suffices to prove uniqueness for the $R_{\mathfrak{p}}$-module $M_{\mathfrak{p}}$. Since the set of all ideals in the valuation domain $R_{\mathfrak{p}}$ form a chain, the uniqueness follows from Proposition 5.22 .

Suppose now that $\mathcal{V}(\mathcal{C})$ is not factorial. By 1. and 2 ., $\mathcal{V}(\mathcal{C}) \cong F \times D$ where $F$ is free abelian and where $D$ is not factorial, but is finitely primary of rank 1 and exponent 1 . Then

$$
\mathrm{c}(F \times D)=\mathrm{t}(F \times D)=\mathrm{c}(D)=\mathrm{t}(D)=2
$$

and hence $F \times D$ is half-factorial.

Remark 5.4. Since a Noetherian Prüfer domain is precisely a Dedekind domain and since in the Noetherian setting the concepts of finitely presented and finitely generated modules coincide, Theorem 5.3 also describes the monoid of all finitely generated modules over a Dedekind domain. Of course, these results can be obtained even more simply from the classical results of Steinitz. In the following section we consider directsum decompositions of yet another class of rings that generalize Dedekind domains.

## 6. Monoids of modules over hereditary Noetherian prime rings

In this final section we study classes $\mathcal{C}$ of finitely generated right modules over hereditary Noetherian prime (HNP) rings, a generalization of Dedekind prime rings (see [43, §5.7]). Module theory over HNP rings is carefully presented in the monograph of Levy and Robson 41, and it is on this work that this section is based. We begin with the arithmetical preparations necessary to state the main result in this
section, Theorem6.5. There we give a characterization of the monoid of stable isomorphism classes of finitely generated projective right modules over HNP rings and use this information to study its arithmetic.

Proposition 6.1. Let $H_{0}$ and $D$ be monoids and define

$$
H=H_{0} \propto D=\left(H_{0} \backslash H_{0}^{\times}\right) \times D \cup H_{0}^{\times} \times\left\{1_{D}\right\}
$$

Then $H$ is a submonoid of $H_{0} \times D$. If $D=\left\{1_{D}\right\}$ or $H_{0}$ is a group, then $H=H_{0}$. Suppose that $D \neq\left\{1_{D}\right\}$ and that $H_{0}$ is not a group.

1. $H^{\times}=H_{0}^{\times} \times\left\{1_{D}\right\}$.
2. $\mathrm{q}(H)=\mathrm{q}\left(H_{0}\right) \times \mathrm{q}(D), H \subset H_{0} \times D$ is not saturated, and $\widehat{H}=\widehat{H}_{0} \times \widehat{D}$. In particular, $H$ is not completely integrally closed and hence not a Krull monoid.
3. The projection $\theta: H \rightarrow H_{0},(a, d) \mapsto a$ is a transfer homomorphism.
4. Suppose $H_{0}$ and $D$ are atomic. Let $u \in \mathcal{A}\left(H_{0}\right)$ and let $d \in D$. Then

$$
\max \left\{\omega\left(H_{0}, u\right), \omega(D, d)\right\} \leq \omega(H,(u, d)) \leq \omega\left(H_{0}, u\right)+\omega(D, d)+\epsilon
$$

where $\epsilon=1$ if $u$ is prime and $d \in D^{\times}$, and $\epsilon=0$ otherwise. In particular, $\omega(H,(u, d))<\infty$ if and only if $\omega\left(H_{0}, u\right)<\infty$ and $\omega(D, d)<\infty$, and if $D$ is not a group, then $\omega(H)=\infty$ and $\mathrm{t}(H)=\infty$.
5. Suppose $H_{0}$ is atomic, $D$ is a group, and let $d \in D$. If $u \in \mathcal{A}\left(H_{0}\right)$ is a prime element, then

$$
\omega(H,(u, d))=2, \tau(H,(u, d))=1 \text { and } \mathrm{t}\left(H,(u, d) H^{\times}\right)=2 .
$$

If $u \in \mathcal{A}\left(H_{0}\right)$ is not a prime element, then

$$
\omega(H,(u, d))=\omega\left(H_{0}, u\right), \tau(H,(u, d))=\tau\left(H_{0}, u\right) \text { and } \mathrm{t}\left(H,(u, d) H^{\times}\right)=\mathrm{t}\left(H_{0}, u H_{0}^{\times}\right) .
$$

In particular,

$$
\mathrm{t}(H)=\max \left\{2, \mathrm{t}\left(H_{0}\right)\right\}, \omega(H)=\max \left\{2, \omega\left(H_{0}\right)\right\} \text { and } \tau(H)=\max \left\{1, \tau\left(H_{0}\right)\right\} .
$$

6. Suppose $H_{0}$ and $D$ are atomic. Let $d \in D$ and $a \in H_{0} \backslash\left(\mathcal{A}\left(H_{0}\right) \cup H_{0}^{\times}\right)$. Then $\mathrm{c}_{H}((a, d))=0$ if

- $D^{\times}=\{1, d\}$, and $\mathrm{Z}_{H_{0}}(a)=\left\{\left(u H_{0}^{\times}\right)^{2}\right\}$ for some $u \in \mathcal{A}\left(H_{0}\right)$, or
- $D$ is reduced, $d=1$ and $\mathrm{c}_{H_{0}}(a)=0$, or
- $D$ is reduced, $d \in \mathcal{A}(D)$ and $Z_{H_{0}}(a)=\left\{\left(u H_{0}^{\times}\right)^{k}\right\}$ for some $u \in \mathcal{A}\left(H_{0}\right)$ and $k \in \mathbb{N}_{\geq 2}$. In any other case, $\mathrm{c}_{H}((a, d))=\max \left\{2, \mathrm{c}_{H_{0}}(a)\right\}$. In particular,

$$
\mathrm{c}(H)=\max \left\{2, \mathrm{c}\left(H_{0}\right)\right\} .
$$

Before proceeding with the proof, the reader may find it useful to note that

$$
H=H_{0} \propto D=\left\{(h, d) \in H_{0} \times D: h \in H_{0}^{\times} \text {only if } d=1_{D}\right\} .
$$

Also, it will be convenient in the proof to introduce the following notation. For a monoid $S$ and elements $a, b \in S$, we write $a \| b$ to denote that $a \mid b$ and $b \nmid a$, that is, $a$ is a strict divisor of $b$.
Proof. It is easily checked that $H$ is a submonoid of $H_{0} \times D$ and that $H=H_{0}$ if $H_{0}=H_{0}^{\times}$or $D=\left\{1_{D}\right\}$. Assume now that $D \neq\left\{1_{D}\right\}$ and $H_{0} \neq H_{0}^{\times}$. Fix $a_{0} \in H_{0} \backslash H_{0}^{\times}$and $d_{0} \in D \backslash\left\{1_{D}\right\}$.

The proof of 1 is clear.
For statement 2, note that we have $\mathrm{q}(H) \subset \mathrm{q}\left(H_{0}\right) \times \mathrm{q}(D)$ and must show the reverse inclusion. Let $a, b \in H_{0}$ and $d, e \in D$. Then $a a_{0}, b a_{0} \in H_{0} \backslash H_{0}^{\times}$and thus $\left(a b^{-1}, d e^{-1}\right)=\left(a a_{0}, d\right)\left(b a_{0}, e\right)^{-1} \in \mathbf{q}(H)$. Therefore $\mathrm{q}\left(H_{0}\right) \times \mathrm{q}(D) \subset \mathrm{q}(H)$. To see that the inclusion $H \subset H_{0} \times D$ is not saturated, note that $\left(1, d_{0}\right) \in \mathrm{q}(H) \cap\left(H_{0} \times D\right)$, but $\left(1, d_{0}\right) \notin H$.

We now show that $\widehat{H}=\widehat{H}_{0} \times \widehat{D}$. Since $H \subset H_{0} \times D$, and since both monoids have the same quotient group, it is certainly true that $\widehat{H} \subset \widehat{H_{0} \times D}=\widehat{H}_{0} \times \widehat{D}$. For the reverse inclusion, let $(x, y) \in \widehat{H}_{0} \times \widehat{D}$. By definition, there exist $c \in H_{0}$ and $d \in D$ such that $(c, d)(x, y)^{n} \in H_{0} \times D$ for all $n \in \mathbb{N}_{0}$. But then $c a_{0} x^{n} \in H_{0} \backslash H_{0}^{\times}$ and $d y^{n} \in D$ for all $n \in \mathbb{N}_{0}$ and thus $\left(c a_{0}, d\right)(x, y)^{n} \in H$. Therefore $(x, y) \in \widehat{H}$. Since $H \neq \widehat{H}, H$ is not completely integrally closed and thus not a Krull monoid.

We now prove statement 3. Clearly $\theta$ is surjective and, by $1, \theta^{-1}\left(H_{0}^{\times}\right)=H^{\times}$. Let $(a, d) \in H$ and $b, c \in H_{0}$ such that $\theta((a, d))=b c$. To establish that $\theta$ is a transfer homomorphism, it will suffice to find $e, f \in D$ such that $(b, e),(c, f) \in H$ and $e f=d$ since, in this case, $(a, d)=(b, e)(c, f), \theta((b, e))=b$ and $\theta((c, f))=c$. If $b$ is not a unit in $H_{0}$, then $e=d$ and $f=1$ gives $(b, d),(c, 1) \in H$ with $e f=d$. Similarly, if $c$ is not a unit, then $e=1$ and $f=d$ gives $(b, 1),(c, d) \in H$ with $e f=d$. If both $b, c \in H_{0}^{\times}$, then $a \in H_{0}^{\times}$and it necessarily follows that $d=1$. Now $e=f=1$ gives $(b, 1),(c, 1) \in H$.

To show 4 we first observe that, for $(a, d),(b, e) \in H,(a, d) \mid(b, e)$ if and only if $a \mid b$ and $d=e$ or $a \| b$ and $d \mid e$. We first prove the correctness of the lower bound and for this we may assume that $k=\omega(H,(u, d))<\infty$. Suppose that $m, k \in \mathbb{N}$ with $k \leq m$ and that $u_{1}, \ldots, u_{m} \in \mathcal{A}(H)$ are such that $u \mid u_{1} \cdot \ldots \cdot u_{m}$. Then $(u, d) \mid\left(u_{1}, d\right) \cdot \ldots \cdot\left(u_{m}, d\right)$, and hence there exists a subproduct of at most $k$ elements that is divisible by $(u, d)$, say $(u, d) \mid\left(u_{1}, d\right) \cdot \ldots \cdot\left(u_{k}, d\right)$. But then $u \mid u_{1} \ldots \ldots u_{k}$ and consequently $\omega\left(H_{0}, u\right) \leq k$. Similarly, suppose $m, k \in \mathbb{N}$ with $k \leq m$ and suppose $d_{1}, \ldots, d_{m} \in \mathcal{A}(D)$ are such that $d \mid d_{1} \cdot \ldots \cdot d_{m}$. Then $(u, d) \mid\left(u^{2}, d_{1}\right) \cdot \ldots \cdot\left(u^{2}, d_{m}\right)$, and so there is a subproduct of at most $k$ elements that is divisible by $(u, d)$, say $(u, d) \mid\left(u^{2}, d_{1}\right) \cdot \ldots \cdot\left(u^{2}, d_{k}\right)$. But then $d \mid d_{1} \cdot \ldots \cdot d_{k}$ and so $\omega(D, d) \leq k$.

We now verify the upper bound, and for this we may assume that $k=\omega(H, u)<\infty$ and $l=\omega(D, d)<\infty$. Suppose that $m, k, l \in \mathbb{N}$ with $m \geq k+l+\epsilon$ and suppose $\left(u_{1}, d_{1}\right), \ldots,\left(u_{m}, d_{m}\right) \in \mathcal{A}(H)$ are such that $(u, d) \mid\left(u_{1}, d_{1}\right) \cdot \ldots \cdot\left(u_{m}, d_{m}\right)$. Then $u \mid u_{1} \cdot \ldots \cdot u_{m}$ and $d \mid d_{1} \cdot \ldots \cdot d_{m}$, whence there are subsets $I, J \subset[1, m]$ with $|I| \leq k$ and $|J| \leq l$ such that $u \mid \prod_{i \in I} u_{i}$ and $d \mid \prod_{j \in J} d_{j}$. Now $|I \cup J| \leq k+l$ and, after renumbering and possibly enlarging $I$ and/or $J$, we have that $I \cup J=[1, k+l]$ with $u \mid u_{1} \cdot \ldots \cdot u_{k+l}$ and $d \mid d_{1} \cdot \ldots \cdot d_{k+l}$. If $u$ is not a prime element, then $k \geq 2$. If $d$ is not a unit then $l \geq 1$. In either of these cases, $k+l \geq 2$ and thus $u \| u_{1} \cdot \ldots \cdot u_{k+l}$. Therefore $(u, d) \mid\left(u_{1}, d_{1}\right) \cdot \ldots \cdot\left(u_{k+l}, d_{k+l}\right)$, and consequently $\omega(H,(u, d)) \leq k+l$. On the other hand, if $u$ is a prime element and $d$ is a unit (i.e., $k+l=1$ ), then $u \| u_{1} \cdot \ldots \cdot u_{k+l+1}$ and hence $(u, d) \mid\left(u_{1}, d_{1}\right) \cdot \ldots \cdot\left(u_{k+l+1}, d_{k+l+1}\right)$.

The claim that $\omega(H,(u, d))<\infty$ if and only if $\omega\left(H_{0}, u\right)<\infty$ and $\omega(D, d)<\infty$ is clear from the nowverified inequalities. If $D$ is not a group, then (since $D$ is atomic) there exists some atom $d \in D$. Then $\omega\left(D, d^{k}\right) \geq k$ for all $k \in \mathbb{N}$, and thus $\omega\left(H,\left(u, d^{k}\right)\right) \geq k$ for all $u \in \mathcal{A}\left(H_{0}\right)$. This implies that $\omega(H)=\infty$ and therefore $\mathrm{t}(H)=\infty$ by Proposition 2.1.

We now prove statement 5 in which case we assume that $D$ is a group. Since, for any atomic monoid $S$ and any non-prime atom $u \in S, \mathrm{t}\left(S, u S^{\times}\right)=\max \{\omega(S, u), \tau(S, u)+1\}$ (see Section 2), it will suffice to establish the claim for the $\omega$ - and $\tau$-invariants.

First observe that since every element of $D$ is a unit, if $(a, e),(b, f) \in H$ with $a \mid b$, then $(a, e) \mid(b, f)(c, g)$ for any $(c, g) \in H \backslash H^{\times}$, and if $a \| b$ then $(a, e) \mid(b, f)$. In particular, if $k \geq 2, u_{1}, \ldots, u_{k} \in \mathcal{A}\left(H_{0}\right)$, and $u \in \mathcal{A}\left(H_{0}\right)$, then $u \mid u_{1} \cdot \ldots \cdot u_{k}$ if and only if $(u, e) \mid\left(u_{1}, e_{1}\right) \cdot \ldots \cdot\left(u_{k}, e_{k}\right)$ for any (equivalently, all) $e, e_{1}, \ldots, e_{k} \in D$.

Since $d \in D^{\times}$, we have that $\omega(D, d)=0$ and thus statement 4 implies that $\omega(H,(u, d))=\omega\left(H_{0}, u\right)$ if $u$ is not a prime element, and $\omega(H,(u, d)) \in[1,2]$ if $u$ is a prime element (since then $\omega\left(H_{0}, u\right)=1$ ). Moreover, $(u, d) \mid\left(u, d d_{0}\right)^{2}$ but $(u, d) \nmid\left(u, d d_{0}\right)$, and hence $\omega(H,(u, d)) \geq 2$ implying that if $u$ is a prime element, then $\omega(H,(u, d))=2$.

A similar argument shows that $\tau(H,(u, d)) \geq 1$ for all $u \in \mathcal{A}\left(H_{0}\right)$. Suppose that $u$ is a prime element and that $k \in \mathbb{N}_{\geq 2},\left(u_{1}, d_{1}\right), \ldots,\left(u_{k}, d_{k}\right) \in \mathcal{A}(H)$, and $(u, d) \mid\left(u_{1}, d_{1}\right) \cdot \ldots \cdot\left(u_{k}, d_{k}\right)$. Then $u \mid u_{1} \cdot \ldots \cdot u_{k}$ and thus $u$ divides one of $u_{1}, \ldots, u_{k}$, say $u \mid u_{1}$. But then $(u, d) \mid\left(u_{1}, d_{1}\right)\left(u_{2}, d_{2}\right)$, showing $\tau(H,(u, d)) \leq 1$.

Recall the definition of the $\tau$-invariant as the supremum of a certain set from equation (1) in Section 2. If $u \in \mathcal{A}\left(H_{0}\right)$ is not a prime element, then the supremum of this set is attained for $k \geq 2$, and since $\mathrm{L}_{H}((a, e))=\mathrm{L}_{H_{0}}(a)$ for all $a \in H_{0}$ and $e \in D$ (by statement 3 ), it is immediate that $\tau(H,(u, \bar{d}))=\tau\left(H_{0}, u\right)$. We note that if $k=2$, then there may be factorizations in $H$ that contribute elements to this set that are not already contributed by factorizations in $H_{0}$. However, if $k=2$, then we necessarily have that $\min \mathrm{L}\left(u^{-1} a\right)=1$ and thus the result is the same as for $k>2$.

We assume now that $H_{0}$ and $D$ are atomic and make use of the transfer homomorphism from statement 3 in order to prove statement 6. By Lemma $2.2,3(\mathrm{~b})$, we have that $\mathrm{c}_{H_{0}}(a) \leq \mathrm{c}_{H}((a, d)) \leq \max \left\{\mathrm{c}(a, \theta), \mathrm{c}_{H_{0}}(a)\right\}$.

We first show that $\mathrm{c}(a, \theta) \leq 2$. For an element $(b, e) \in H$ we write $\overline{(b, e)}=(b, e) H^{\times}$for its class in $H_{\text {red }}$. Suppose that $z=\overline{\left(u_{1}, d_{1}\right)} \cdot \ldots \cdot \overline{\left(u_{k}, d_{k}\right)}$ and $z^{\prime}=\overline{\left(u_{1}, d_{1}^{\prime}\right)} \cdot \ldots \cdot \overline{\left(u_{k}, d_{k}^{\prime}\right)}$ with $k, k^{\prime} \geq 2, u_{1}, \ldots, u_{k} \in$ $\underline{\mathcal{A}\left(H_{0}\right)}$, and $d_{1}, \ldots, d_{k}, d_{1}^{\prime}, \ldots, d_{k}^{\prime} \in D$ are two factorizations of $(a, d)$ lying in the same fibre of $\theta$. Then $\overline{\left(u_{1}, d_{1}\right)} \cdot \ldots \cdot \overline{\left(u_{k-1}, d_{k-1} d_{k}\right)} \overline{\left(u_{k}, 1\right)}$ is also a factorization of $(a, d)$ lying in the same fibre and

$$
\mathrm{d}\left(\overline{\left(u_{1}, d_{1}\right)} \cdot \ldots \cdot \overline{\left(u_{k-1}, d_{k-1}\right)} \overline{\left(u_{k}, d_{k}\right)}, \overline{\left(u_{1}, d_{1}\right)} \cdot \ldots \cdot \overline{\left(u_{k-1}, d_{k-1} d_{k}\right)} \overline{\left(u_{k}, 1\right)}\right) \leq 2 .
$$

Inductively, we find a 2 -chain from $z$ to $\overline{\left(u_{1}, d_{1} \cdot \ldots \cdot d_{k}\right)} \overline{\left(u_{2}, 1\right)} \cdot \ldots \cdot \overline{\left(u_{k}, 1\right)}$. Similarly, we find a 2 -chain from $z^{\prime}$ to $\overline{\left(u_{1}, d_{1}^{\prime} \cdot \ldots \cdot d_{k}^{\prime}\right)} \overline{\left(u_{2}, 1\right)} \cdot \ldots \cdot \overline{\left(u_{k}, 1\right)}$ and hence a 2 -chain from $z$ to $z^{\prime}$. This shows that $\mathrm{c}(a, \theta) \leq 2$.

Recall that $\mathrm{c}_{H}((a, d))=0$ is equivalent to $(a, d)$ having a unique factorization in $H$ and, in any other case, $\mathrm{c}_{H}((a, d)) \geq 2$. Therefore, to establish the remaining claims, it will suffice to show that in each of the cases listed, the element $(a, d)$ has a unique factorization and that in any other case, $(a, d)$ has at least two distinct factorizations. From the inequalities already proven, we then obtain that $\mathrm{c}_{H}((a, d))=\max \left\{2, \mathrm{c}_{H_{0}}(a)\right\}$.

We first suppose that $D$ is not reduced. Assume that there exists $\varepsilon \in D^{\times} \backslash\{1, d\}$ and let $k \in \mathbb{N}_{\geq 2}$ and $u_{1}, \ldots, u_{k} \in \mathcal{A}\left(H_{0}\right)$ with $a=u_{1} \cdot \ldots \cdot u_{k}$. Then

$$
\overline{\left(u_{1}, d\right)} \overline{\left(u_{2}, 1\right)} \cdot \ldots \cdot \overline{\left(u_{k}, 1\right)} \quad \text { and } \overline{\left(u_{1}, d \varepsilon^{-1}\right)} \overline{\left(u_{2}, \varepsilon\right)} \overline{\left(u_{3}, 1\right)} \cdot \ldots \cdot \overline{\left(u_{k}, 1\right)}
$$

are two distinct factorizations of $(a, d)$ since $\left(d \varepsilon^{-1}, \varepsilon\right) \neq\{(d, 1),(1, d)\}$. Therefore, in this case, $\mathrm{c}_{H}((a, d)) \geq 2$.
We now consider the case where $D^{\times}=\{1, d\}$ with $d \neq 1$. If $Z_{H_{0}}(a)=\left\{\left(u H_{0}^{\times}\right)^{2}\right\}$ for some $u \in \mathcal{A}\left(H_{0}\right)$, then the unique factorization of $(a, d)$ is $\overline{(u, d)} \overline{(u, 1)}$. The remaining cases are as follows.

- If $a$ has two distinct factorizations, then so does $(a, d)$.
- If $a$ has a unique factorization represented by $u_{1} \cdot \ldots \cdot u_{k}$ for some $k \in \mathbb{N}_{\geq 2}$ and $u_{1}, \ldots, u_{k} \in \mathcal{A}\left(H_{0}\right)$ with $u_{1} \nsucceq u_{2}$, then

$$
\overline{\left(u_{1}, d\right)} \overline{\left(u_{2}, 1\right)} \overline{\left(u_{3}, 1\right)} \cdot \ldots \overline{\left(u_{k}, 1\right)} \quad \text { and } \overline{\left(u_{1}, 1\right)} \overline{\left(u_{2}, d\right)} \overline{\left(u_{3}, 1\right)} \cdot \ldots \cdot \overline{\left(u_{k}, 1\right)}
$$

are two distinct factorizations of $(a, d)$.

- If $a \simeq u^{k}$ for some $k \in \mathbb{N}_{\geq 3}$ and $u \in \mathcal{A}\left(H_{0}\right)$, then $\overline{(u, d)} \overline{(u, 1)}^{k-1}$ and $\overline{(u, d)}^{3} \overline{(u, 1)}^{k-3}$ are two distinct factorizations of $(a, d)$ (here we use $\left.d^{2}=1\right)$.
Now suppose that $D$ is reduced. If $d=1$ and, for some $k \in \mathbb{N}_{\geq 2}$ and $u_{1}, \ldots, u_{k} \in \mathcal{A}\left(H_{0}\right), a=u_{1} \cdot \ldots \cdot u_{k}$ represents the unique factorization of $a$ in $H_{0}$, then $\overline{\left(u_{1}, 1\right)} \cdot \ldots \cdot \overline{\left(u_{k}, 1\right)}$ is the unique factorization of $(a, d)$ in $H$. If $d \in \mathcal{A}(D)$ and for some $u \in \mathcal{A}\left(H_{0}\right)$ and $k \in \mathbb{N}_{\geq 2}, a \simeq u^{k}$ represents the unique factorization of $a$ in $H_{0}$, then $\overline{(u, d)} \overline{(u, 1)}^{k-1}$ is the unique factorization of $(a, d)$ in $H$.

We now show that in each of the remaining cases there exist two distinct factorizations of $(a, d)$ in $H$. This is clear if $a$ has two distinct factorizations in $H_{0}$, and we may assume in the following that $a$ has a unique factorization represented by $a=u_{1} \cdot \ldots \cdot u_{k}$ for some $k \in \mathbb{N}_{\geq 2}$ and $u_{1}, \ldots, u_{k} \in \mathcal{A}\left(H_{0}\right)$.

- If $d \neq 1$ is not an atom, there must exist $d_{1}, d_{2} \in D \backslash\{1\}$ such that $d=d_{1} d_{2}$. Then $(a, d)=$ $\left(u_{1}, d\right)\left(u_{2}, 1\right) \cdot \ldots \cdot\left(u_{k}, 1\right)$ and $(a, d)=\left(u_{1}, d_{1}\right)\left(u_{2}, d_{2}\right)\left(u_{3}, 1\right) \cdot \ldots \cdot\left(u_{k}, 1\right)$ give rise to two distinct factorizations of $(a, d)$.
- If $d \in \mathcal{A}(D)$ and there exist two distinct factors in the factorization of $a$, say $u_{1} \nsucceq u_{2}$, then $(a, d)=$ $\left(u_{1}, d\right)\left(u_{2}, 1\right)\left(u_{3}, 1\right) \cdot \ldots \cdot\left(u_{k}, 1\right)$ and $(a, d)=\left(u_{1}, 1\right)\left(u_{2}, d\right)\left(u_{3}, 1\right) \cdot \ldots \cdot\left(u_{k}, 1\right)$ give rise to two distinct factorizations of $(a, d)$.
Note that we have so far showed that $\mathrm{c}\left(H_{0}\right) \leq \mathrm{c}(H) \leq \max \left\{2, \mathrm{c}\left(H_{0}\right)\right\}$. We now show that $2 \leq \mathrm{c}(H)$. Since $H_{0}$ is atomic, but not a group, there exists $a \in H_{0} \backslash H_{0}^{\times}$. Taking a power of $a$ if necessary, we may further assume that $a \in H_{0} \backslash \mathcal{A}\left(H_{0}\right)$. If $D$ is not reduced, then $2 \leq \mathrm{c}((a, 1)) \leq \mathrm{c}(H)$. If $D$ is reduced, then, since $D \neq\left\{1_{D}\right\}$ by assumption, there exists $d \in D \backslash D^{\times}$and, again taking a power of $a$ if necessary, we may assume $d \in D \backslash \mathcal{A}(D)$. Thus $2 \leq \mathrm{c}((a, d)) \leq \mathrm{c}(H)$.

We note that if $H_{0}=\mathbb{N}_{0}$ and $D$ is a group, then $H=\mathbb{N}_{0} \propto D$ as above is a finitely primary monoid of rank 1 and exponent 1 (as discussed at the beginning of Section 5). Also, if $H_{0}, D$ and $E$ are monoids,
then $\left(H_{0} \propto D\right) \propto E=H_{0} \propto(D \times E)$. Before defining, in Definition 6.3, a monoid that will later be used to describe the monoid of stable isomorphism classes of finitely generated projective right modules over an HNP ring, we give a simple lemma that will be used in the proof of Proposition 6.4

Lemma 6.2. Let $G$ be an additive abelian group and let $G_{0} \subset G$ be a subset.

1. If $G_{0}=\{g\}$ for some torsion element $g \in G$, then $\mathcal{A}\left(G_{0}\right)=\left\{g^{\operatorname{ord}(g)}\right\}$.
2. Let $n \in \mathbb{N}_{\geq 2}$, let $\left(e_{1}, \ldots, e_{n-1}\right)$ be a family of independent elements of $G$ each of infinite order, and let

$$
G_{0}=\left\{A_{i} e_{i}: i \in[1, n-1]\right\} \cup\left\{-\sum_{i=1}^{n-1} B_{i} e_{i}\right\}
$$

where $A_{i}, B_{i} \in \mathbb{N}$ with $\operatorname{gcd}\left(A_{i}, B_{i}\right)=1$ for each $i \in[1, n-1]$. Then $\mathcal{A}\left(G_{0}\right)=\{U\}$ where

$$
U=\prod_{i=1}^{n-1}\left(A_{i} e_{i}\right)^{\frac{B_{i} L}{A_{i}}}\left(-\sum_{i=1}^{n-1} B_{i} e_{i}\right)^{L} \quad \text { with } \quad L=\operatorname{lcm}\left\{A_{i}: i \in[1, n-1]\right\}
$$

In particular, $\mathcal{B}\left(G_{0}\right)$ is factorial in each case.
Proof. If $\mathcal{A}\left(G_{0}\right)=\{S\}$, then $S$ is a prime element, in which case $\mathcal{B}\left(G_{0}\right)$ is factorial.
Statement 1 is clear and we now prove statement 2. From the definition of $L$, we have $A_{i} \mid B_{i} L$. Moreover, $\sigma(U)=\sum_{i=1}^{n-1}\left(\frac{B_{i} L}{A_{i}} A_{i}-B_{i} L\right) e_{i}=0 \in G$ which shows that $U \in \mathcal{B}\left(G_{0}\right)$. Suppose that $S=\prod_{i=1}^{n-1}\left(A_{i} e_{i}\right)^{k_{i}}(-$ $\left.\sum_{i=1}^{n-1} B_{i} e_{i}\right)^{k} \in \mathcal{B}\left(G_{0}\right)$ for some $k_{1}, \ldots, k_{n}, k \in \mathbb{N}_{0}$. Note that the $k_{1}, \ldots, k_{n}$ are uniquely determined by $k$ and $\sigma(S)=0$. Therefore, to establish that $U$ is an atom of $\mathcal{B}\left(G_{0}\right)$, and in fact the unique atom, it will suffice to show that $L \mid k$. Since $\sigma(S)=0$, we have that $A_{i} k_{i}=k B_{i}$ for each $i \in[1, n-1]$. Since $\operatorname{gcd}\left(A_{i}, B_{i}\right)=1$, this implies $A_{i} \mid k$ and hence $L=\operatorname{lcm}\left\{A_{i}: i \in[1, n-1]\right\} \mid k$.

Definition 6.3. Let $\Omega$ be a set containing a designated element 0 and let $\mathbf{c} \in \mathbb{Q}_{>0}^{\Omega}$ such that $c_{0}=1$ and $c_{i} \in \mathbb{N}$ for all but finitely many $i \in \Omega$. Define

$$
\mathbb{N}_{0}^{\Omega}(\mathbf{c})=\left\{\mathbf{x} \in \mathbb{N}_{0}^{\Omega}: x_{0}>0 \text { and }\left|\operatorname{supp}_{\mathbb{Z}^{\Omega}}\left(\mathbf{x}-x_{0} \mathbf{c}\right)\right|<\infty\right\} \cup\{\mathbf{0}\},
$$

that is, $\mathbb{N}_{0}^{\Omega}(\mathbf{c})$ consists of those vectors which are almost everywhere equal to a non-zero multiple (determined by the coordinate $x_{0}$ ) of $\mathbf{c}$, together with the vector $\mathbf{0}$. We write $\ell(\mathbf{x})=x_{0}$. If $|\Omega|<\infty$, then $\mathbb{N}_{0}^{\Omega}(\mathbf{c}) \cong$ $\left(\mathbb{N} \times \mathbb{N}_{0}^{\Omega \backslash\{0\}}\right) \cup\{\mathbf{0}\} \cong\left(\mathbb{N} \times \mathbb{N}_{0}^{(\Omega \backslash\{0\})}\right) \cup\{\mathbf{0}\}$. Let $\Lambda$ be a (possibly empty) set of finite, pairwise disjoint subsets of $\Omega \backslash\{0\}$ each containing at least two elements and, for each $I \in \Lambda$, let $C_{I}=\sum_{i \in I} c_{i}$. We assume that $C_{I} \in \mathbb{N}$ for all $I \in \Lambda$. Define

$$
\mathbb{N}_{0}^{\Omega}(\mathbf{c}, \Lambda)=\left\{\mathbf{x} \in \mathbb{N}_{0}^{\Omega}(\mathbf{c}): \sum_{i \in I} x_{i}=C_{I} \ell(\mathbf{x}) \text { for all } I \in \Lambda\right\}
$$

By definition, $\mathbb{N}_{0}^{\Omega}(\mathbf{c})$ and $\mathbb{N}_{0}^{\Omega}(\mathbf{c}, \Lambda)$ are reduced submonoids of $\left(\mathbb{N}_{0}^{\Omega},+\right)$. However, note that the inclusion $\mathbb{N}_{0}^{\Omega}(\mathbf{c}) \subset \mathbb{N}_{0}^{\Omega}$ is not saturated. Indeed, if $\mathbf{x}, \mathbf{y} \in \mathbb{N}_{0}^{\Omega}(\mathbf{c})$, then $\mathbf{x}$ divides $\mathbf{y}$ if and only if $\mathbf{y}-\mathbf{x} \in \mathbb{N}_{0}^{\Omega}(\mathbf{c})$, that is, $\mathbf{x} \leq \mathbf{y}$ and either $x_{0}<y_{0}$ or $\mathbf{x}=\mathbf{y}$.

Proposition 6.4. Let $H=\mathbb{N}_{0}^{\Omega}(\mathbf{c}, \Lambda)$ with $\Omega$, $\mathbf{c}$ and $\Lambda$ as in Definition 6.3. Then $H$ is a saturated submonoid of $\mathbb{N}_{0}^{\Omega}(\mathbf{c})$ and the map $\ell: H \rightarrow\left(\mathbb{N}_{0},+\right)$ is a transfer homomorphism. In particular, $\mathbf{x} \in H$ is an atom if and only if $\ell(\mathbf{x})=1$, and $H$ is half-factorial.

1. If $|\Omega|=1$, then $H=\mathbb{N}_{0}$.
2. Suppose $2 \leq|\Omega|<\infty$ and $\bigcup \Lambda=\Omega \backslash\{0\}$. Write $\Omega=[0, r]$ with $r \in \mathbb{N}$, and $\Lambda=\left\{I_{1}, \ldots, I_{n}\right\}$ with $n \in \mathbb{N}$ and $\Omega \backslash\{0\}=I_{1} \uplus \cdots \uplus I_{n}$. Set $C_{i}=C_{I_{i}}$ for all $i \in[1, n]$.
(a) The map $j$ : $H \hookrightarrow \mathbb{N}_{0}^{r}, \mathbf{x} \mapsto\left(x_{1}, \ldots, x_{r}\right)$ is a divisor theory (note that $x_{0}$ is omitted), and hence $H$ is a Krull monoid. Then $G=\mathbb{Z}^{r} / \mathbf{q}(j(H))$ is the divisor class group and $G_{P}=\left\{\mathbf{e}_{i}+\mathbf{q}(j(H)): i \in\right.$ $[1, r]\} \subset G$ is the set of classes containing prime divisors. There is a monomorphism

$$
\varphi^{*}: G \rightarrow\left\{\begin{array}{ll}
\mathbb{Z} / C_{1} \mathbb{Z} & \text { if } n=1, \\
\mathbb{Z}^{n-1} & \text { if } n \geq 2
\end{array} \quad\right. \text { with the following properties: }
$$

If $n=1$, then $\varphi^{*}$ is an isomorphism, $\varphi^{*}\left(G_{P}\right)=\left\{1+C_{1} \mathbb{Z}\right\}$, and the unique class in $G_{P}$ contains precisely $\left|I_{1}\right|$ prime divisors. Let $n \geq 2$ and, for all $i \in[1, n-1]$, set $A_{i}=\frac{C_{n}}{\operatorname{gcd}\left(C_{i}, C_{n}\right)}$ and $B_{i}=\frac{C_{i}}{\operatorname{gcd}\left(C_{i}, C_{n}\right)}$. Then $\varphi^{*}(G)$ is a full rank subgroup of $\mathbb{Z}^{n-1}$ and

$$
\varphi^{*}\left(G_{P}\right)=\left\{A_{i} \mathbf{e}_{i}: i \in[1, n-1]\right\} \cup\left\{-\sum_{i=1}^{n-1} B_{i} \mathbf{e}_{i}\right\}
$$

For each $i \in[1, n-1]$, the class mapped onto the element $A_{i} \mathbf{e}_{i}$ contains precisely $\left|I_{i}\right|$ prime divisors and the class mapped onto $-\sum_{i=1}^{n-1} B_{i} \mathbf{e}_{i}$ contains precisely $\left|I_{n}\right|$ prime divisors.
(b) An element $\mathbf{u} \in H$ is an atom if and only if $\sum_{i \in I} u_{i}=C_{I}$ for all (equivalently, any) $I \in \Lambda$.
(c) We have that $\mathrm{c}(H) \leq 2$. Moreover, the monoid $H$ is factorial if and only if $|\Lambda|=1$ and $C_{1}=1$. In this case, the prime elements are precisely the $\mathbf{e}_{0}+\mathbf{e}_{i}$ for $i \in[1, r]$. If $H$ is not factorial, then $H$ contains no prime elements.
(d) For any atom $\mathbf{u} \in \mathcal{A}(H)$ we have

$$
\sum_{I \in \Lambda}\left(C_{I}-\min _{i \in I} u_{i}\right) \leq \omega(H, \mathbf{u}) \leq \sum_{I \in \Lambda} C_{I}
$$

(e) If $H$ is not factorial, then $\mathrm{t}(H)=\omega(H)=\sum_{I \in \Lambda} C_{I}$.
3. If $2 \leq|\Omega|<\infty$ but $\bigcup \Lambda \subsetneq \Omega \backslash\{0\}$, set $\Omega^{\prime}=\bigcup \Lambda \cup\{0\}$. Let $\mathbf{c}^{\prime} \in \mathbb{Q}_{>0}^{\Omega^{\prime}}$ be defined by $c_{i}^{\prime}=c_{i}$ for all $i \in \Omega^{\prime}$ and let $H^{\prime}=\mathbb{N}_{0}^{\Omega^{\prime}}\left(\mathbf{c}^{\prime}, \Lambda\right)$. Then $H \cong H^{\prime} \propto \mathbb{N}_{0}^{|\Omega|-\left|\Omega^{\prime}\right|}$. For each $(\mathbf{u}, \mathbf{x}) \in \mathcal{A}(H)$,

$$
\max \left\{\omega\left(H^{\prime}, \mathbf{u}\right),|\mathbf{x}|\right\} \leq \omega(H,(\mathbf{u}, \mathbf{x})) \leq \omega\left(H^{\prime}, \mathbf{u}\right)+|\mathbf{x}|+1
$$

and thus $\omega(H,(\mathbf{u}, \mathbf{x}))<\infty$ and $\mathrm{t}(H,(\mathbf{u}, \mathbf{x}))<\infty$. Moreover, $H$ is a half-factorial FF -monoid with $\mathrm{c}(H)=2$ and $\mathrm{t}(H)=\infty$. Also, $H$ is not a Krull monoid.
4. If $|\Omega|=\infty$, then
(a) $\left|\mathrm{Z}_{H}(\mathbf{x})\right|=\infty$ for all $\mathbf{x} \in H \backslash\left(\mathcal{A}(H) \cup H^{\times}\right)$. In particular, $H$ is not a FF -monoid and is therefore not a submonoid of a free abelian monoid,
(b) $\mathrm{c}_{H}(\mathbf{x})=2$ for all $\mathbf{x} \in H \backslash\left(\mathcal{A}(H) \cup H^{\times}\right)$and
(c) $\omega(H, \mathbf{u})=\tau(H, \mathbf{u})=\mathrm{t}(H, \mathbf{u})=\infty$ for all $\mathbf{u} \in \mathcal{A}(H)$.

Proof. Suppose that $\mathbf{x}, \mathbf{y} \in H$ such that $\mathbf{x}$ divides $\mathbf{y}$ in $\mathbb{N}_{0}^{\Omega}(\mathbf{c})$. That is, $\mathbf{y}-\mathbf{x} \in \mathbb{N}_{0}^{\Omega}(\mathbf{c})$. For each $I \in \Lambda$, we have $\sum_{i \in I} y_{i}-x_{i}=C_{I}(\ell(\mathbf{y})-\ell(\mathbf{x}))=C_{I} \ell(\mathbf{y}-\mathbf{x})$, and thus $\mathbf{y}-\mathbf{x} \in H$. Therefore the inclusion $H \subset \mathbb{N}_{0}^{\Omega}(\mathbf{c})$ is saturated. We now show that $\ell: H \rightarrow \mathbb{N}_{0}$ is a transfer homomorphism. Clearly this map is surjective and, since $H$ is reduced, it will suffice to show: If $\mathbf{x} \in H$ and $\ell(\mathbf{x})=k+l$ with $k, l \in \mathbb{N}_{0}$, then there exist $\mathbf{y}, \mathbf{z} \in H$ such that $\mathbf{x}=\mathbf{y}+\mathbf{z}$ and $\ell(\mathbf{y})=k, \ell(\mathbf{z})=l$.

If one of $k$ or $l$ is 0 we may, without restriction, assume that $l=0$. Then $\mathbf{y}=\mathbf{x}$ and $\mathbf{z}=\mathbf{0}$ give the result. From now on we assume that $k, l>0$. Let $\Lambda^{\prime} \subset \Lambda$ consist of those finitely many $I \in \Lambda$ with $I \cap \operatorname{supp}(\mathbf{x}-\ell(\mathbf{x}) \mathbf{c}) \neq \emptyset$. For $i \in \bigcup \Lambda^{\prime}$, let $y_{i}^{\prime}, z_{i}^{\prime} \in \mathbb{N}_{0}$ be such that $y_{i}^{\prime}+z_{i}^{\prime}=x_{i}$ and such that, for all $I \in \Lambda^{\prime}$, it holds that $\sum_{i \in I} y_{i}^{\prime}=C_{I} k$ and $\sum_{i \in I} z_{i}^{\prime}=C_{I} l$. Define $\mathbf{y}$ and $\mathbf{z} \in \mathbb{N}_{0}^{\Omega}$ by

$$
y_{i}=\left\{\begin{array}{ll}
y_{i}^{\prime} & \text { if } i \in \bigcup \Lambda^{\prime}, \\
k c_{i} & \text { if } i \in \Omega \backslash \bigcup \Lambda^{\prime},
\end{array} \quad z_{i}= \begin{cases}z_{i}^{\prime} & \text { if } i \in \bigcup \Lambda^{\prime} \\
l c_{i} & \text { if } i \in \Omega \backslash \bigcup \Lambda^{\prime}\end{cases}\right.
$$

Then $\mathbf{y}, \mathbf{z} \in \mathbb{N}_{0}^{\Omega}(\mathbf{c}, \Lambda)$ with $\ell(\mathbf{y})=k, \ell(\mathbf{z})=l$, and $\mathbf{x}=\mathbf{y}+\mathbf{z}$ as required.
Statement 1 is clear and thus we suppose that $2 \leq|\Omega|<\infty$ and $\Omega \backslash\{0\}=\bigcup \Lambda$.

We first verify statement $2(\mathrm{a})$. Since $\sum_{i \in I_{1}} x_{i}=C_{1} x_{0}$ and $C_{1} \neq 0, j$ is injective. We now check that $j$ is a divisor homomorphism. Let $\mathbf{x}, \mathbf{y} \in H$ with $j(\mathbf{x}) \leq j(\mathbf{y})$. If $\mathbf{x}=\mathbf{y}$, there is nothing to show. If $\mathbf{x} \neq \mathbf{y}$, then there necessarily exists $I \in \Lambda$ with $\sum_{i \in I} x_{i}<\sum_{i \in I} y_{i}$ and hence $x_{0}<y_{0}$. Consequently, $\mathbf{x}$ divides $\mathbf{y}$ in $H$ and thus $j$ is a divisor homomorphism. To prove that $j$ is a divisor theory, we need to show that each standard basis vector $\mathbf{e}_{i} \in \mathbb{N}_{0}^{r}$ is the greatest common divisor of a finite, non-empty set in the image of $j$. Let $i \in[1, r]$, let $I_{0} \in \Lambda$ be such that $i \in I_{0}$, and let $i^{\prime} \in I_{0} \backslash\{i\}$. Define

$$
\begin{aligned}
& \mathbf{x}=\mathbf{e}_{0}+\mathbf{e}_{i}+\left(C_{I_{0}}-1\right) \mathbf{e}_{i^{\prime}}+\sum_{I \in \Lambda \backslash\left\{I_{0}\right\}} C_{I} \mathbf{e}_{\min I}, \text { and } \\
& \mathbf{y}=\mathbf{e}_{0}+C_{I_{0}} \mathbf{e}_{i}+\sum_{I \in \Lambda \backslash\left\{I_{0}\right\}} C_{I} \mathbf{e}_{\max I}
\end{aligned}
$$

Then $\mathbf{x}, \mathbf{y} \in H$ and, recalling that $|I| \geq 2$ for all $I \in \Lambda, \operatorname{gcd}_{\left(\mathbb{N}_{0}^{r},+\right)}(j(\mathbf{x}), j(\mathbf{y}))=\mathbf{e}_{i}$. Thus $j: H \hookrightarrow \mathbb{N}_{0}^{r}$ is a divisor theory, $H$ is a Krull monoid, and the class group of $H$ is $\mathrm{q}\left(\mathbb{N}_{0}^{r}\right) / \mathrm{q}(j(H))=\mathbb{Z}^{r} / \mathrm{q}(j(H))$.

We now determine the structure of the divisor class group of $H$ and determine the set of classes containing prime divisors. Suppose that $n=1$ and define $\varphi^{*}: \mathbb{Z}^{r} \rightarrow \mathbb{Z} / C_{1} \mathbb{Z}$ by $\mathbf{x} \mapsto \sum_{i=1}^{r} x_{i}+C_{1} \mathbb{Z}$. We see immediately that $\varphi^{*}$ is an epimorphism and $\operatorname{ker}\left(\varphi^{*}\right)=\mathrm{q}(j(H))$ follows easily. Therefore $G \cong \mathbb{Z} / C_{1} \mathbb{Z}$. Since $\varphi^{*}\left(\mathbf{e}_{i}\right)=$ $1+C_{1} \mathbb{Z}$ for all $i \in[1, r], G_{P}$ is as claimed. Now suppose that $n \geq 2$ and let $\varphi^{*}: \mathbb{Z}^{r} \rightarrow \mathbb{Z}^{n-1}$ be defined by

$$
\mathbf{x} \mapsto\left(A_{1} \sum_{i \in I_{1}} x_{i}-B_{1} \sum_{i \in I_{n}} x_{i}, \quad A_{2} \sum_{i \in I_{2}} x_{i}-B_{2} \sum_{i \in I_{n}} x_{i}, \ldots, \quad A_{n-1} \sum_{i \in I_{n-1}} x_{i}-B_{n-1} \sum_{i \in I_{n}} x_{i}\right) .
$$

One easily checks that $\operatorname{ker}\left(\varphi^{*}\right)=\mathrm{q}(j(H))$, showing that $G$ embeds into $\mathbb{Z}^{n-1}$ via $\varphi^{*}$. Considering the images of $\mathbf{e}_{1}, \ldots, \mathbf{e}_{r}$ under $\varphi^{*}$, the description of $G_{P}$ given in statement $2(\mathrm{a})$ follows, and we see that $\varphi^{*}(G)$ is a subgroup of full rank of $\mathbb{Z}^{n-1}$.

For statement $2(\mathrm{~b})$, note that an element $\mathbf{u} \in H$ is an atom if and only if $\ell(\mathbf{u})=1$, and this is the case if and only if $\sum_{i \in I} u_{i}=C_{I}$ for all (equivalently, any) $I \in \Lambda$.

Consider statement 2(c). By Lemma Lemma 6.2, $\mathcal{B}\left(G_{P}\right)$ is factorial and hence $\mathrm{c}\left(G_{P}\right)=0$. Thus Proposition 2.3 implies that $\mathrm{c}(H) \leq \max \left\{2, \mathrm{c}\left(G_{P}\right)\right\} \leq 2$. If $|\Lambda|=1$ and $C_{1}=1$, then $j$ is surjective. Thus $H \cong \mathbb{N}_{0}^{r}$, showing that $H$ is factorial. The prime elements of $\left(\mathbb{N}_{0}^{r},+\right)$ are simply the standard basis vectors $\mathbf{e}_{i}$ for $i \in[1, r]$, and their preimages under $j$ are precisely the elements $\mathbf{e}_{0}+\mathbf{e}_{i} \in H$ for $i \in[1, r]$. If $|\Lambda|>1$ or $C_{1}>1$, then no atom is prime and thus $H$ is not factorial. Let $\mathbf{u} \in \mathcal{A}(H)$. The lower bound given in $2(\mathrm{~d})$, which we will soon verify, implies that $\omega(H, \mathbf{u}) \geq 2$ unless $|\Lambda|=1, C_{1}=2$, and $u_{1}=u_{2}=1$. Note that $\min _{i \in I} u_{i} \leq\left\lfloor C_{I} / 2\right\rfloor$ for each $I \in \Lambda$ and thus, in this case, $\mathbf{u}$ is not a prime element. In the remaining case, one easily checks that again, $\mathbf{u}$ is not a prime element.

We now verify the bounds on the omega invariant as given in $2(\mathrm{~d})$. If $|\Lambda|=1$ and $C_{1}=1$, then $H$ is factorial and the inequalities hold trivially. We assume from now on that this is not the case and hence $\sum_{I \in \Lambda} C_{I} \geq 2$. We first show that $\omega(H, \mathbf{u}) \leq \sum_{I \in \Lambda} C_{I}$. Let $k \in \mathbb{N}_{\geq 2}$ and let $\mathbf{v}_{1}, \ldots, \mathbf{v}_{k} \in \mathcal{A}(H)$ be such that $\mathbf{u}$ divides $\sum_{i=1}^{k} \mathbf{v}_{i}$. If $J \subset[1, k]$, then $\mathbf{u}$ divides $\sum_{j \in J} \mathbf{v}_{j}$ if and only if $\mathbf{u} \leq \sum_{j \in J} \mathbf{v}_{j}$. Since $\mathbf{u} \leq \sum_{i=1}^{k} \mathbf{v}_{i}$ and $\sum_{i \in \Omega \backslash\{0\}} u_{i}=\sum_{I \in \Lambda} C_{I}$, we can recursively construct a subset $J \subset[1, k]$ of size at most $\sum_{I \in \Lambda} C_{I}$ such that $\mathbf{u} \leq \sum_{j \in J} \mathbf{v}_{j}$. This is done by adding, in each step, a vector $\mathbf{v}_{j}$ with $v_{j, i}>0$ for some $i \in \Omega \backslash\{0\}$ for which $u_{i}<\sum_{j \in J} v_{j, i}$.

We now prove the lower bound. By renumbering the coordinates if necessary, we may assume $u_{\min I}=$ $\min _{i \in I} u_{i}$ for all $I \in \Lambda$. For $i \in \Omega$, let $I \in \Lambda$ be the unique set containing $i$ and define

$$
\mathbf{v}_{i}=\mathbf{e}_{0}+\mathbf{e}_{i}+\left(C_{I}-1\right) \mathbf{e}_{\min I}+\sum_{I^{\prime} \in \Lambda \backslash\{I\}} C_{I^{\prime}} \mathbf{e}_{\min } I^{\prime} \in \mathcal{A}(H) .
$$

Since $\mathbf{u}=\sum_{i \in \Omega} u_{i} \mathbf{e}_{i}=\sum_{I \in \Lambda} \sum_{i \in I} u_{i} \mathbf{e}_{i}$, it is clear that $\mathbf{u}$ divides $\sum_{I \in \Lambda} \sum_{i \in I \backslash\{\min I\}} u_{i} \mathbf{v}_{i}$. For $I \in \Lambda$ with $C_{I}>1$, this is clear from the definition of the $\mathbf{v}_{i}$ 's. If $C_{I}=1$, then $u_{\min I}=0$ since $|I| \geq 2$, yet $\mathbf{u}$ does not
divide any proper subsum. Therefore

$$
\omega(H, \mathbf{u}) \geq \sum_{I \in \Lambda}\left[\left(\sum_{i \in I} u_{i}\right)-\min _{i \in I} u_{i}\right]=\sum_{I \in \Lambda}\left(C_{I}-\min _{i \in I} u_{i}\right) .
$$

We now prove statement 2(e). By definition, $\omega(H)=\sup _{\mathbf{u} \in \mathcal{A}(H)} \omega(H, \mathbf{u})$. Since $H$ is half-factorial, $\mathrm{t}(H)=\omega(H)$. The element

$$
\mathbf{u}=\mathbf{e}_{0}+\sum_{I \in \Lambda} C_{I} \mathbf{e}_{\min I} \in H
$$

is an atom and the bounds in $2(\mathrm{~d})$ imply that $\omega(H, \mathbf{u})=\sum_{I \in \Lambda} C_{I}$. Therefore $\omega(H) \geq \sum_{I \in \Lambda} C_{I}$ and the upper bound again follows from 2 (d).

Consider statement 3 and let $D=\mathbb{N}_{0}^{|\Omega|-\left|\Omega_{0}\right|}$. The isomorphism $H \cong H^{\prime} \propto D$ is immediate from the definitions. Because $D$ is factorial, we have $\omega(D, \mathbf{x})=|\mathbf{x}|$ for all $\mathbf{x} \in D$. Proposition 6.16 together with 2(c) implies that $\mathrm{c}(H)=2$, and Proposition 6.14 implies the remaining inequalities. As a submonoid of a free abelian monoid, $H$ is an FF-monoid ([28, Corollary 1.5.7]). However, since it is not completely integrally closed in its quotient group (by Proposition 6.12 2), it cannot be a Krull monoid.

We suppose for the remainder of the proof that $|\Omega|=\infty$. For notational ease we assume $\mathbb{N}_{0} \subset \Omega$.
Consider statement $4(\mathrm{a})$. To show $\left|\mathrm{Z}_{H}(\mathbf{x})\right|=\infty$ it suffices to show that there are infinitely many atoms of $H$ dividing $\mathbf{x}$. By definition of $H$, there exists a finite subset $\Omega^{\prime} \subset \Omega$ having the property that $I \cap \Omega^{\prime} \neq \emptyset$ already implies $I \subset \Omega^{\prime}$ for all $I \in \Lambda$ and such that for each $i \in \Omega \backslash \Omega^{\prime}$ it holds that $x_{i}=\ell(x) c_{i}$ and $c_{i} \in \mathbb{N}$. We may assume that $\mathbb{N}_{0} \subset \Omega \backslash \Omega^{\prime}$. Moreover, we may assume that for each $I \in \Lambda$ we have $|I \cap \mathbb{N}| \leq 1$ and, if $i \in I \cap \mathbb{N}$, then $c_{i}=\min _{i^{\prime} \in I} c_{i^{\prime}}$.

For each $i \in \mathbb{N}$, define $\mathbf{u}_{i} \in \mathbb{N}_{0}^{\Omega}$ as follows: If there exists $I \in \Lambda$ with $i \in I$, then let $i^{\prime} \in I \backslash\{i\}$. Note also that $i^{\prime} \in \Omega \backslash \Omega^{\prime}$ by the choice of $\Omega^{\prime}$. We set $u_{i, i}=0, u_{i, i^{\prime}}=c_{i}+c_{i^{\prime}}$, and $u_{i, j}=c_{j}$ for each $j \in \Omega \backslash\left(\Omega^{\prime} \cup\left\{i, i^{\prime}\right\}\right)$. On the other hand, if $i$ is not contained in any $I \in \Lambda$, we set $u_{i, i}=0$ and $u_{i, j}=c_{j}$ for each $j \in \Omega \backslash\left(\Omega^{\prime} \cup\{i\}\right)$. In either case, we can choose $u_{i, j} \leq x_{i, j}$ for each $j \in \Omega^{\prime}$ such that $\sum_{j \in I} u_{i, j}=C_{I}$ for all $I \in \Lambda$ with $I \subset \Omega^{\prime}$. Therefore $\sum_{j \in I} u_{i, j}=C_{I}$ for all $I \in \Lambda$ and $u_{i, j}=c_{j}$ for all but finitely many $j \in \Omega$. Therefore $\mathbf{u}_{i} \in \mathcal{A}(H)$. By construction, $\mathbf{u}_{i} \leq \mathbf{x}$. In the first case, due to the minimal choice of $c_{i}$, we have $u_{i, i^{\prime}}=c_{i}+c_{i^{\prime}} \leq 2 c_{i^{\prime}} \leq \ell(\mathbf{x}) c_{i^{\prime}}=x_{i^{\prime}}$, with the last equality holding since $i^{\prime} \in \Omega \backslash \Omega^{\prime}$. Since $\ell\left(\mathbf{u}_{i}\right)=1<\ell(\mathbf{x})$, $\mathbf{u}_{i}$ divides $\mathbf{x}$. Clearly these atoms are pairwise distinct. By what we have just shown, $H$ is not an FF-monoid and therefore cannot a submonoid of a free abelian monoid by [28, Corollary 1.5.7].

We now compute the catenary degree of $H$ as stated in 4(b). From 4(a), every nonzero, non-atom element has at least two distinct factorizations and hence $\mathrm{c}_{H}(\mathbf{x}) \geq 2$ for all $\mathbf{x} \in H$. We now show that $\mathrm{c}_{H}(\mathbf{x}) \leq 2$ by projecting to the finite case. Suppose that $\mathbf{x}=\mathbf{u}_{1}+\cdots+\mathbf{u}_{k}=\mathbf{v}_{1}+\cdots+\mathbf{v}_{k}$ for some $k \in \mathbb{N}_{\geq 2}$ and that $\mathbf{u}_{1}, \ldots, \mathbf{u}_{k}, \mathbf{v}_{1}, \ldots, \mathbf{v}_{k} \in \mathcal{A}(H)$. Let $\Omega^{\prime} \subset \Omega$ be the smallest subset containing

$$
\left\{j \in \Omega: u_{i, j} \neq c_{j} \text { or } v_{i, j} \neq c_{j} \text { for some } i \in[1, k]\right\} \cup\{0\}
$$

and having the property that whenever $I \cap \Omega^{\prime} \neq \emptyset$ for $I \in \Lambda$, then already $I \subset \Omega^{\prime}$. Observe that $\Omega^{\prime}$ is finite and that $\mathbf{u}_{i, j}=\mathbf{v}_{i^{\prime}, j}=c_{j}$ for all $i, i^{\prime} \in[1, k]$ and $j \in \Omega \backslash \Omega^{\prime}$. Set $\Lambda^{\prime}=\left\{I \in \Lambda: I \subset \Omega^{\prime}\right\}$ and let $\mathbf{c}^{\prime} \in \mathbb{N}_{0}^{\Omega^{\prime}}$ be defined by $c_{i}^{\prime}=c_{i}$ for all $i \in \Omega^{\prime}$. Define $H^{\prime}=\mathbb{N}_{0}^{\Omega^{\prime}}\left(\mathbf{c}^{\prime}, \Lambda^{\prime}\right)$ and note that there is a canonical projection $\pi: H \rightarrow H^{\prime}$. By the finiteness of $\Omega^{\prime}$ we immediately have that $\mathrm{c}\left(H^{\prime}\right) \leq 2$. Thus there exists a sequence of factorizations $z_{1}^{\prime}, \ldots, z_{l}^{\prime} \in \mathrm{Z}\left(H^{\prime}\right)$ of $\pi(\mathbf{x})$ with $z_{1}^{\prime}=\pi\left(\mathbf{u}_{1}\right) \cdot \ldots \cdot \pi\left(\mathbf{u}_{k}\right), z_{l}^{\prime}=\pi\left(\mathbf{v}_{1}\right) \cdot \ldots \cdot \pi\left(\mathbf{v}_{k}\right)$ and such that $\mathrm{d}\left(z_{i}^{\prime}, z_{i+1}^{\prime}\right) \leq 2$ for all $i \in[1, l-1]$. For any factor $\mathbf{w}^{\prime}$ occurring in some $z_{j}^{\prime}$, we lift it to $\mathbf{w} \in H$ by setting $w_{j}=u_{1, j}$ for all $j \in \Omega \backslash \Omega^{\prime}$. Then $\mathbf{w} \in \mathcal{A}(H)$ and we can lift the factorizations $z_{1}^{\prime}, \ldots, z_{l}^{\prime}$ of $\pi(\mathbf{x})$ to factorizations $z_{1}, \ldots, z_{l}$ of $\mathbf{x}$. These factorizations form a sequence connecting the two factorization of $\mathbf{x}$ that we began with and have the property that $\mathrm{d}\left(z_{i}, z_{i+1}\right) \leq 2$ for all $i \in[1, l-1]$.

Finally, we consider the infinitude of the invariants given in $4(\mathrm{c})$. Let $\mathbf{u} \in \mathcal{A}(H)$ and $k \in \mathbb{N}_{\geq 2}$. We will show that $\omega(H, \mathbf{u}) \geq k$, as the other invariants are then equal to $\omega(H, u)$ by half-factoriality. Note that $\mathbf{u}$ is necessarily non-zero in all but finitely many coordinates and so we may assume $u_{i} \neq 0$ for $i \in[1, k]$.

Moreover, since $\Omega$ is infinite and each $I \in \Lambda$ is finite, we may suppose that no two $i, j \in[1, k]$ belong to the same set $I \in \Lambda$. (For the following argument it would suffice that $I \not \subset[1, k]$.) Modifying finitely many coordinates of $\mathbf{u}$ as needed, we can therefore construct $\mathbf{u}_{1}, \ldots, \mathbf{u}_{k} \in \mathcal{A}(H)$ such that $u_{i, j}=u_{j} \delta_{i, j}$ for all $i, j \in[1, k]$, and $u_{i, j} \geq u_{j}$ for all $i \in[1, k]$ and $j \in \Omega \backslash[1, k]$. Then $\mathbf{u}$ divides $\mathbf{u}_{1}+\cdots+\mathbf{u}_{k}$, but for all $j \in[1, k]$, $u_{j}>u_{1, j}+\cdots+u_{k, j}-u_{j, j}=0$, showing that $\mathbf{u}$ divides no proper subsum.

Now with the appropriate arithmetical results, we apply Propositions 6.1 and 6.4 to study direct-sum decompositions of modules over HNP rings. For a hereditary Noetherian prime ring (HNP ring) $R$, we consider the class $\mathcal{C}$ of finitely generated right $R$-modules. In the noncommutative setting there are two invariants (the genus and the Steinitz class) which describe the stable isomorphism class of a finitely generated projective right $R$-module $P$. In general, however, the isomorphism class of $P$ is determined by its stable isomorphism class only if $\operatorname{udim}(P) \geq 2$ and such a module is indecomposable if and only if $\operatorname{udim}(P)=1$. Thus the forthcoming description of the direct-sum decomposition of finitely generated projective right $R$-modules is one where the indecomposable factors are determined up to stable isomorphism. If $R$ has the additional property that any two finitely generated projective right $R$-modules that are stably isomorphic are already isomorphic, then this result is a description up to isomorphism.

Let $V$ and $W$ be two simple right $R$-modules. Then $W$ is called a successor of $V$ and $V$ is called a predecessor of $W$ if $\operatorname{Ext}_{R}^{1}(V, W) \neq \mathbf{0}$. Let $\mathcal{W}$ be a set of representatives (of isomorphism classes) of the simple unfaithful right $R$-modules. We note that every $V \in \mathcal{W}$ is contained in a unique tower of $R([41, \S 19])$. A tower $\mathcal{T}$ is a finite set of simple right $R$-modules, ordered with respect to the successor relationship, and having the following structure: Every tower $\mathcal{T}$ is either cyclically ordered and each simple module in $\mathcal{T}$ is unfaithful, in which case we say that $\mathcal{T}$ is a cycle tower, or $\mathcal{T}$ is linearly ordered and only the first module (the only module in the tower not having a predecessor) is faithful, in which case $\mathcal{T}$ is said to be a faithful tower. The length of a tower is the number of distinct modules contained in it, and a tower is non-trivial if it contains more than one module.

We briefly recall the notions of rank, genus and the Steinitz class of a finitely generated projective right $R$-module, as these invariants are used to describe stable isomorphism classes of finitely generated projective right $R$-modules. We refer the reader to [41, $\S 33$ and $\S 35]$ for additional details. Let $P$ be a finitely generated projective right $R$-module and let $V \in \mathcal{W}$. Then $M=\operatorname{ann}(V)$ is a maximal ideal of $R, R / M$ is simple artinian, and the rank of $M$ at $V$, denoted by $\rho(P, V) \in \mathbb{N}_{0}$, is defined to be the length of the $R / M$-module $P / P M$. If $\mathcal{T}$ is a cycle tower, we set $\rho(P, \mathcal{T})=\sum_{V \in \mathcal{T}} \rho(P, V)$. Let $\operatorname{modspec}(R)$ denote the set of isomorphism classes of all unfaithful simple right $R$-modules together with the trivial module $\mathbf{0}$. For a finitely generated projective right $R$-module $P$, we set $\Psi_{V}(P)=\rho(P, V)$ if $V$ is an unfaithful simple right $R$-module, and $\Psi_{\mathbf{0}}(P)=\operatorname{udim}(P)$. Then

$$
\Psi(P)=\left(\Psi_{V}(P)\right)_{V \in \operatorname{modspec}(R)} \in \mathbb{N}_{0}^{\operatorname{modspec}(R)}
$$

is called the genus of $P$. Two finitely generated projective right $R$-modules $P$ and $Q$ are stably isomorphic if there exists a finitely generated projective right $R$-module $X$ such that $P \oplus X \cong Q \oplus X$. We denote by $[P]$ the stable isomorphism class of $P$. The direct sum operation on modules induces the structure of a commutative semigroup on the set of stable isomorphism classes, and by $\mathbf{K}_{0}(R)$ we denote its quotient group. The genus $\Psi$ induces a homomorphism $\Psi^{+}: \mathbf{K}_{0}(R) \rightarrow \mathbb{Z}^{\operatorname{modspec}(R)}$, and $G(R)=\operatorname{Ker}\left(\Psi^{+}\right)$is called the ideal class group of $R$. By choosing a base point set $\mathcal{B}$ of non-zero finitely generated projective right $R$-modules consisting of exactly one module in each genus, and such that $\mathcal{B}$ is closed, up to isomorphism, under direct sums, we can associate to any nonzero finitely generated projective right $R$-module $P$ a class $\mathcal{S}(P)=[P]-[B] \in G(R)$, where $B \in \mathcal{B}$. We call $\mathcal{S}(P)$ the Steinitz class of $P$ and we set $\mathcal{S}(\mathbf{0})=\mathbf{0}$.

We are now ready to characterize the semigroup of stable isomorphism classes of finitely generated projective right modules over an HNP ring. We note that the characterization of factoriality in Theorem 6.5 was already obtained by Levy and Robson in [41, Theorem 39.5].

Theorem 6.5. Let $R$ be an HNP ring and let $H$ be the semigroup of stable isomorphism classes of finitely generated projective right $R$-modules with operation induced by the direct sum of modules.

1. Let $\Omega$ denote the set of isomorphism classes of all unfaithful simple right $R$-modules which are contained in a non-trivial tower, together with the trivial module, denoted by $\mathbf{0}$. For $V \in \Omega \backslash\{\mathbf{0}\}$, let $c_{V}=$ $\frac{\rho(R, V)}{\operatorname{udim}(R)} \in \mathbb{Q}_{>0}$ and $c_{\mathbf{0}}=1$. Finally, let $\Lambda$ be the set of all non-trivial cycle towers. Then

$$
H \cong \mathbb{N}_{0}^{\Omega}(\mathbf{c}, \Lambda) \propto G(R)
$$

with the product $\propto$ as defined in Proposition 6.1. In particular, $H$ is half-factorial, $\mathrm{c}(H) \leq 2$ and the following are equivalent:
(a) $H$ is factorial (i.e., $R$ satisfies stable uniqueness).
(b) $G(R)=\mathbf{0}$ and either $R$ has no non-trivial towers (i.e., $R$ is a Dedekind prime ring), or $R$ has a unique non-trivial tower $\mathcal{T}$ which is a cycle tower and $\rho(R, \mathcal{T})=\operatorname{udim}(R)$.
2. If $R$ has infinitely many non-trivial towers, then $\mathrm{t}(H,[U])=\omega(H,[U])=\infty$ for each indecomposable finitely generated projective right $R$-module $U$.
3. Suppose that $R$ has only finitely many non-trivial towers. Then $\mathrm{t}(H,[U]), \omega(H,[U])<\infty$ for each indecomposable finitely generated projective right $R$-module $U$. If $R$ has at least one non-trivial faithful tower, then $\mathrm{t}(H)=\omega(H)=\infty$. If $R$ has no non-trivial faithful tower, only finitely many non-trivial cycle towers $\mathcal{T}_{1}, \ldots, \mathcal{T}_{n}$ with $n \in \mathbb{N}_{0}$, and $H$ is not factorial, then

$$
\mathrm{t}(H)=\omega(H)=\frac{1}{\operatorname{udim}(R)} \sum_{i=1}^{n} \rho\left(R, \mathcal{T}_{i}\right)
$$

Proof. It will suffice to establish that $H \cong \mathbb{N}_{0}^{\Omega}(\mathbf{c}, \Lambda) \propto G(R)$ as the remaining claims will then follow from Proposition 6.4. The genus and the Steinitz class are both additive on direct sums and thus give rise to a monoid homomorphism $H \rightarrow \mathbb{N}_{0}^{\operatorname{modspec}(R)} \propto G(R),[P] \mapsto(\Psi(P), \mathcal{S}(P))$. We note that $\mathcal{S}(\mathbf{0})=\mathbf{0}$ and that $\mathbf{0}$ is the only module of uniform dimension zero. The genus satisfies a number of necessary conditions, namely that $\Psi_{V}(P)=c_{V} \Psi_{0}(P)$ for almost all $V \in \operatorname{modspec}(R)$ and that it has standard rank at every cycle tower, that is, $\rho(P, \mathcal{T})=\sum_{V \in \mathcal{T}} c_{V} \Psi_{0}(P)$ for all cycle towers $\mathcal{T}$. In particular, if $V$ is contained in a trivial tower, then this is necessary a cycle tower and hence implies $\Psi_{V}(P)=c_{V} \Psi_{0}(P)$. Thus instead of $\Psi(P)$ we may consider $\Psi^{\prime}(P) \subset \mathbb{N}_{0}^{\Omega}$ were we omit the components corresponding to unfaithful simple right $R$-modules that are contained in a trivial tower. We obtain a homomorphism

$$
\Phi: H \rightarrow \mathbb{N}_{0}^{\Omega}(\mathbf{c}, \Lambda) \propto G(R),[P] \mapsto\left(\Psi^{\prime}(P), \mathcal{S}(P)\right)
$$

The main theorem of Levy and Robson ([41, Theorem 35.13]) implies that the genus and the Steinitz class are independent invariants, and that up to the stated conditions on the rank, all values can be obtained. In other words, $\Phi$ is an isomorphism.

Remark 6.6. 1. If $R$ is such that each two stably isomorphic right $R$-modules are isomorphic, then $H=$ $\mathcal{V}\left(\mathcal{C}_{\text {proj }}\right)$. If this is not the case, then $H$ still provides information on direct sum decompositions of finitely generated projective modules with the summands determined up to stable isomorphism. Specifically, we have the following: Clearly, if $P=U_{1} \oplus \cdots \oplus U_{k}$ for some $k \in \mathbb{N}$ and indecomposable right $R$-modules $U_{1}, \ldots, U_{k}$, then $[P]=\left[U_{1}\right]+\cdots+\left[U_{k}\right]$ is a factorization of $[P]$ in $H$. Let $P$ be a right $R$-module with $\operatorname{udim}(P) \geq 2$ (that is, $P$ is neither the zero module nor indecomposable). If $[P]=\left[U_{1}\right]+\cdots+\left[U_{k}\right]$ for some $k \in \mathbb{N}_{\geq 2}$ and atoms $\left[U_{1}\right], \ldots,\left[U_{k}\right]$ of $H$, then $[P]=\left[U_{1} \oplus \cdots \oplus U_{k}\right]$ and, since $\operatorname{udim}(P) \geq 2$, 41, Theorem 34.6] implies that $P \cong U_{1} \oplus \cdots \oplus U_{k}$. Therefore, a factorization of $[P]$ in $H$ gives rise to one of $P$ into indecomposables, with the stable isomorphism classes of the indecomposable summands determined by the factorization of $[P]$ in $H$.
2. An HNP ring $R$ has finitely many non-trivial towers if and only if it is a multichain idealizer from a Dedekind prime ring $S$ (41, Proposition 30.5]). A sufficient (but not necessary) condition for there to be no non-trivial faithful towers is for $R$ to be right (equivalently, left) bounded ([41, Lemma 18.2]).
3. Let $R$ be an HNP ring and let $\mathcal{C}$ denote the class of finitely generated right $R$-modules. Then $\mathcal{V}(\mathcal{C}) \cong$ $\mathcal{V}\left(\mathcal{C}_{\text {tor }}\right) \times \mathcal{V}\left(\mathcal{C}_{\text {proj }}\right)$ and $\mathcal{V}\left(\mathcal{C}_{\text {tor }}\right)$ is factorial.
Proof. By [41, Corollary 12.16(ii), Theorem 12.18], every finitely generated right $R$-module decomposes uniquely as a direct sum of a torsion module and a torsion-free module. The first has finite length and hence has a unique decomposition into indecomposables, while the second is projective.

A detailed description of $\mathcal{V}\left(\mathcal{C}_{\text {tor }}\right)$ is given in [41, §41].

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