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# NORM INEQUALITIES RELATED TO THE MATRIX GEOMETRIC MEAN OF NEGATIVE POWER 

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#### Abstract

In this paper, we show norm inequalities related to the matrix geometric mean of negative power for positive definite matrices: For positive definite matrices $A$ and $B$, $$
\left\|e^{(1-\beta) \log A+\beta \log B}\right\|\|\leq\| A \natural_{\beta} B\|\leq\| A^{1-\beta} B^{\beta} \|
$$ for every unitarily invariant norm and $-1 \leq \beta \leq-\frac{1}{2}$, where the $\beta$-quasi geometric mean $A \mathfrak{\natural}_{\beta} B$ is defined by $A \mathfrak{\natural}_{\beta} B=A^{\frac{1}{2}}\left(A^{-\frac{1}{2}} B A^{-\frac{1}{2}}\right)^{\beta} A^{\frac{1}{2}}$. For our purposes, we show the Ando-Hiai log-majorization of negative power.


1 Introduction. Let $\mathbb{M}_{n}=\mathbb{M}_{n}(\mathbb{C})$ be the algebra of $n \times n$ complex matrices and denote the matrix absolute value of any $A \in \mathbb{M}_{n}$ by $|A|=\left(A^{*} A\right)^{\frac{1}{2}}$. For $A \in \mathbb{M}_{n}$, we write $A \geq 0$ if $A$ is positive semidefinite and $A>0$ if $A$ is positive definite, that is, $A$ is positive and invertible. For two Hermitian matrices $A$ and $B$, we write $A \geq B$ if $A-B \geq 0$, and it is called the Löwner ordering. A norm $\|\cdot\| \|$ on $\mathbb{M}_{n}$ is said to be unitarily invariant if $\|U X V\|=\|X\|$ for all $X \in \mathbb{M}_{n}$ and unitary $U, V$.

Let $A$ and $B$ be two positive definite matrices. The arithmetic-geometric mean inequality says that

$$
\begin{equation*}
A \not \sharp_{\alpha} B \leq(1-\alpha) A+\alpha B \quad \text { for all } \alpha \in[0,1] \text {, } \tag{1.1}
\end{equation*}
$$

where the $\alpha$-geometric mean $A \sharp_{\alpha} B$ is defined by

$$
A \sharp_{\alpha} B=A^{\frac{1}{2}}\left(A^{-\frac{1}{2}} B A^{-\frac{1}{2}}\right)^{\alpha} A^{\frac{1}{2}} \quad \text { for all } \alpha \in[0,1] \text {, }
$$

also see [11]. As another matrix geometric mean, we recall that the chaotic geometric mean $A \diamond_{\alpha} B$ is defined by

$$
A \diamond_{\alpha} B=e^{(1-\alpha) \log A+\alpha \log B} \quad \text { for all } \alpha \in \mathbb{R}
$$

also see [5, Section 3.5]. If $A$ and $B$ commute, then $A \diamond_{\alpha} B=A \sharp_{\alpha} B=A^{1-\alpha} B^{\alpha}$ for $\alpha \in[0,1]$. In [4], Bhatia and Grover showed precise norm estimations of the arithmeticgeometric mean inequality (1.1) as follows: For each $\alpha \in[0,1]$ and any unitarily invariant norm $\|\cdot\|$

$$
\begin{aligned}
\left\|A \sharp_{\alpha} B\right\| & \leq\left\|A \diamond_{\alpha} B\right\| \leq\left\|B^{\frac{\alpha}{2}} A^{1-\alpha} B^{\frac{\alpha}{2}}\right\| \\
& \leq\left\|\frac{1}{2}\left(A^{1-\alpha} B^{\alpha}+B^{\alpha} A^{1-\alpha}\right)\right\| \leq\left\|A^{1-\alpha} B^{\alpha}\right\| \leq\|(1-\alpha) A+\alpha B\|
\end{aligned}
$$

[^0]and
\[

$$
\begin{aligned}
\left\|A \not \sharp_{\alpha} B\right\| & \leq\left\|A \diamond_{\alpha} B\right\| \\
& \leq\left\|\left(B^{\frac{\alpha p}{2}} A^{(1-\alpha) p} B^{\frac{\alpha p}{2}}\right)^{\frac{1}{p}}\right\| \leq\left\|\left((1-\alpha) A^{p}+\alpha B^{p}\right)^{\frac{1}{p}}\right\| \quad \text { for all } p>0 .
\end{aligned}
$$
\]

For convenience in symbolic expression, we define $A \natural_{\beta} B$ for $\beta \in[-1,0)$ and positive definite matrices $A, B$ as follows:

$$
\begin{equation*}
A \mathfrak{h}_{\beta} B=A^{\frac{1}{2}}\left(A^{-\frac{1}{2}} B A^{-\frac{1}{2}}\right)^{\beta} A^{\frac{1}{2}} \quad \text { for all } \beta \in[-1,0), \tag{1.2}
\end{equation*}
$$

whose formula is the same as $\sharp_{\alpha}$. Though $A \natural_{\beta} B$ for $\beta \in[-1,0)$ are not matrix means in the sense of Kubo-Ando theory [11], it is known in [7] that $A \natural_{\beta} B$ have matrix mean like properties for any positive definite matrices $A$ and $B$. Thus we call (1.2) the $\beta$-quasi geometric mean for $\beta \in[-1,0)$. For more detail, see [7].

On the other hand, the following reverse arithmetic-geometric mean inequality holds:

$$
(1-\beta) A+\beta B \leq A \natural_{\beta} B \quad \text { for all } \beta \in[-1,0),
$$

also see [8]. Though we have no relation among $A$ Ł $_{\beta} B, A \diamond_{\beta} B$ and $A^{1-\beta} B^{\beta}$ for $\beta \in[-1,0)$ under the Löwner ordering, it follows from a proof similar to Bhatia-Grover's one in [4] that for each $\beta \in \mathbb{R}$ and any unitarily invariant norm $\|\cdot\|$

$$
\begin{equation*}
\left\|A \diamond_{\beta} B\right\| \leq\left\|B^{\frac{\beta}{2}} A^{1-\beta} B^{\frac{\beta}{2}}\right\| \leq\left\|\frac{1}{2}\left(A^{1-\beta} B^{\beta}+B^{\beta} A^{1-\beta}\right)\right\| \leq\left\|A^{1-\beta} B^{\beta}\right\| . \tag{1.3}
\end{equation*}
$$

Also, by the Lie-Trotter formula $\lim _{t \rightarrow 0}\left(e^{\frac{t}{2} B} e^{t A} e^{\frac{t}{2} B}\right)^{\frac{1}{t}}=e^{A+B}$ and the Araki-Cordes inequality $\left\|B^{t} A^{t} B^{t}\right\| \leq\left\|(B A B)^{t}\right\|$ for all $t \in[0,1]$, also see [3, Exercise IX.1.5,Theorem IX.2.10], it follows that for each $\beta \in[-1,0)$

$$
\begin{equation*}
\left\|A \diamond_{\beta} B\right\| \leq\left\|\left(B^{\frac{\beta q}{2}} A^{(1-\beta) q} B^{\frac{\beta q}{2}}\right)^{\frac{1}{q}}\right\| \tag{1.4}
\end{equation*}
$$

holds for all $q>0$ and $\left\|\left(B^{\frac{\beta q}{2}} A^{(1-\beta) q} B^{\frac{\beta q}{2}}\right)^{\frac{1}{q}}\right\|$ decreases to $\left\|A \diamond_{\beta} B\right\|$ as $q \downarrow 0$. It is natural to ask what is the estimate of the $\beta$-quasi geometric mean in the norm inequalities (1.3) and (1.4) for $\beta \in[-1,0)$.

In this paper, we show norm inequalities related to the $\beta$-quasi geometric mean of negative power, the chaotic geometric mean $A \diamond_{\beta} B$ and $A^{1-\beta} B^{\beta}$ for positive definite matrices $A, B$. Moreover, we show precise norm estimations of the reverse arithmeticgeometric mean inequality under the assumption $A \geq B$. For our purposes, we need the Ando-Hiai log-majorization of negative power.

2 Preliminaries. In this section, we have some preliminary results on the log majorization of matrices. For Hermitian matrices $H, K$ the weak majorization $H \prec_{w} K$ means that

$$
\sum_{i=1}^{k} \lambda_{i}(H) \leq \sum_{i=1}^{k} \lambda_{i}(K) \quad \text { for } k=1,2, \ldots, n,
$$

where $\lambda_{1}(H) \geq \cdots \geq \lambda_{n}(H)$ and $\lambda_{1}(K) \geq \cdots \geq \lambda_{n}(K)$ are the eigenvalues of $H$ and $K$ respectively. Further, the majorization $H \prec K$ means that $H \prec_{w} K$ and the equality holds

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for $k=n$ in the above, i.e., $\operatorname{Tr} H=\operatorname{Tr} K$. For $A, B \geq 0$ let us write $A \prec_{w(\log )} B$ and refer to the weak $\log$ majorization if

$$
\prod_{i=1}^{k} \lambda_{i}(A) \leq \prod_{i=1}^{k} \lambda_{i}(B) \quad \text { for } k=1,2, \ldots, n
$$

Further the $\log$ majorization $A \prec_{(\log )} B$ means that $A \prec_{w(\log )} B$ and the equality holds for $k=n$ in the above, i.e.,

$$
\prod_{i=1}^{n} \lambda_{i}(A)=\prod_{i=1}^{n} \lambda_{i}(B) \quad \text { i.e., } \quad \operatorname{det} A=\operatorname{det} B
$$

Note that when $A, B>0$ the $\log$ majorization $A \prec_{(\log )} B$ is equivalent to $\log A \prec \log B$. It is known that for positive semidefinite $A, B \geq 0$,

$$
A \prec_{w(\log )} B \Longrightarrow A \prec_{w} B \Longrightarrow\|A\| \leq\|B\|
$$

for any unitarily invariant norm. See $[1,12]$ for theory of majorization for matrices.
For each matrix $X$ and $k=1,2, \ldots, n$, let $C_{k}(X)$ denote the $k$-fold antisymmetric tensor power of $X$. See [12] for details. Then (1)-(3) below are basic facts, (4) is easily seen from (2) and (3), and (5) follows from the Binet-Cauchy theorem.

Lemma 2.1. (1) $C_{k}\left(X^{*}\right)=C_{k}(X)^{*}$.
(2) $C_{k}(X Y)=C_{k}(X) C_{k}(Y)$ for every pair of matrices $X, Y$.
(3) $C_{k}\left(X^{-1}\right)=C_{k}(X)^{-1}$ for nonsingular $X$.
(4) $C_{k}\left(A^{p}\right)=C_{k}(A)^{p}$ for every positive definite $A>0$ and all $p \in \mathbb{R} \backslash\{0\}$.
(5) For every $A>0, \prod_{i=1}^{k} \lambda_{i}(A)=\lambda_{1}\left(C_{k}(A)\right)$ for $k=1,2, \ldots, n$ and consequently, for $A, B>0, \lambda_{1}\left(C_{k}(A)\right) \leq \lambda_{1}\left(C_{k}(B)\right)$ for all $k=1, \ldots, n$ if and only if $A \prec_{w(\log )} B$.

3 Ando-Hiai Log-Majorization of negative power. For $0 \leq \alpha \leq 1$, the matrix $\alpha$ geometric mean is the matrix mean corresponding to the matrix monotone function $t^{\alpha}$. Note that $A \sharp_{\alpha} B=B \sharp_{1-\alpha} A$ and if $A B=B A$ then $A \sharp_{\alpha} B=A^{1-\alpha} B^{\alpha}$, and $(A, B) \mapsto A \sharp_{\alpha} B$ is jointly monotone, also see [5, Lemma 3.2].

On the other hand, the $\beta$-quasi geometric mean for $\beta \in[-1,0)$ has the following properties in [7]; for any positive definite matrices $A, B$ and $C$
(i) consistency with scalars: If $A$ and $B$ commute, then $A \natural_{\beta} B=A^{1-\beta} B^{\beta}$.
(ii) homogeneity: $(\alpha A) দ_{\beta}(\alpha B)=\alpha\left(A \natural_{\beta} B\right)$ for all $\alpha>0$.
(iii) right reverse monotonicity: $B \leq C$ implies $A \natural_{\beta} B \geq A \natural_{\beta} C$.

We recall the log-majorization theorem due to Ando-Hiai [2]: For each $\alpha \in[0,1]$

$$
A^{r} \not \sharp_{\alpha} B^{r} \prec_{(\log )}\left(A \not \sharp_{\alpha} B\right)^{r} \quad \text { for } r \geq 1 \text {, }
$$

or equivalently

$$
\left(A^{p} \sharp_{\alpha} B^{p}\right)^{\frac{1}{p}} \prec_{(\log )}\left(A^{q} \sharp_{\alpha} B^{q}\right)^{\frac{1}{q}} \quad \text { for } 0<q<p .
$$

To show the main theorem related to the $\beta$-quasi geometric mean for $\beta \in[-1,0)$, we need the following Ando-Hiai log-majorization of negative power $\beta \in[-1,0)$ :

Theorem 3.1. For every positive definite matrices $A, B>0$ and $\beta \in[-1,0)$,

$$
\begin{equation*}
A^{r} \natural_{\beta} B^{r} \prec_{(\log )}\left(A \natural_{\beta} B\right)^{r} \quad \text { for all } 0<r \leq 1 \tag{3.1}
\end{equation*}
$$

or equivalently

$$
\begin{align*}
\left(A দ_{\beta} B\right)^{r} \prec_{(\log )}\left(A^{r} \vdash_{\beta} B^{r}\right) & \text { for all } r \geq 1,  \tag{3.2}\\
\left(A^{q} দ_{\beta} B^{q}\right)^{\frac{1}{q}} \prec_{(\log )}\left(A^{p} দ_{\beta} B^{p}\right)^{\frac{1}{p}} & \text { for all } 0<q \leq p . \tag{3.3}
\end{align*}
$$

Proof. The equivalence of (3.1)-(3.3) is immediate. It is easy to see by Lemma 2.1 that for $k=1, \ldots, n$

$$
C_{k}\left(A^{r} \natural_{\beta} B^{r}\right)=C_{k}(A)^{r} \natural_{\beta} C_{k}(B)^{r}
$$

and

$$
C_{k}\left(\left(A \natural_{\beta} B\right)^{r}\right)=\left(C_{k}(A) \natural_{\beta} C_{k}(B)\right)^{r}
$$

Also,

$$
\operatorname{det}\left(A^{r} \natural_{\beta} B^{r}\right)=(\operatorname{det} A)^{r(1-\beta)}(\operatorname{det} B)^{r \beta}=\operatorname{det}\left(A \natural_{\beta} B\right)^{r} .
$$

Hence, in order to prove (3.1), it suffices to show that

$$
\begin{equation*}
\lambda_{1}\left(A^{r} দ_{\beta} B^{r}\right) \leq \lambda_{1}\left(A \natural_{\beta} B\right)^{r} \quad \text { for all } 0<r \leq 1 \tag{3.4}
\end{equation*}
$$

For this purpose we may prove that $A \natural_{\beta} B \leq I$ implies $A^{r} \natural_{\beta} B^{r} \leq I$, because both sides of (3.4) have the same order of homogeneity for $A, B$, so that we can multiply $A, B$ by a positive constant.

First let us assume $\frac{1}{2} \leq r \leq 1$ and write $r=1-\varepsilon$ with $0 \leq \varepsilon \leq \frac{1}{2}$. Let $C=A^{\frac{1}{2}} B^{-1} A^{\frac{1}{2}}$. Then $B^{-1}=A^{-\frac{1}{2}} C A^{-\frac{1}{2}}$ and $A \natural_{\beta} B=A^{\frac{1}{2}} C^{-\beta} A^{\frac{1}{2}}$. If $A \natural_{\beta} B \leq I$, then $C^{-\beta} \leq A^{-1}$ so that $A \leq C^{\beta}$ and $A^{\varepsilon} \leq C^{\beta \varepsilon}$ for $0 \leq \varepsilon \leq \frac{1}{2}$ by Löwner-Heinz inequality. Since $-\beta \in(0,1]$ and $1-\varepsilon \in\left[\frac{1}{2}, 1\right]$, we now get

$$
\begin{aligned}
A^{r} \mathfrak{h}_{\beta} B^{r} & =A^{\frac{1-\varepsilon}{2}}\left(A^{\frac{\varepsilon-1}{2}} B^{1-\varepsilon} A^{\frac{\varepsilon-1}{2}}\right)^{\beta} A^{\frac{1-\varepsilon}{2}} \\
& =A^{\frac{1-\varepsilon}{2}}\left(A^{\frac{1-\varepsilon}{2}}\left(B^{-1}\right)^{1-\varepsilon} A^{\frac{1-\varepsilon}{2}}\right)^{-\beta} A^{\frac{1-\varepsilon}{2}} \\
& =A^{\frac{1-\varepsilon}{2}}\left(A^{\frac{1-\varepsilon}{2}}\left(A^{-\frac{1}{2}} C A^{-\frac{1}{2}}\right)^{1-\varepsilon} A^{\frac{1-\varepsilon}{2}}\right)^{-\beta} A^{\frac{1-\varepsilon}{2}} \\
& =A^{\frac{1-\varepsilon}{2}}\left(A^{-\frac{\varepsilon}{2}}\left[A \sharp_{1-\varepsilon} C\right] A^{-\frac{\varepsilon}{2}}\right)^{-\beta} A^{\frac{1-\varepsilon}{2}} \\
& =A^{\frac{1}{2}-\varepsilon}\left[A^{\varepsilon} \sharp_{-\beta}\left(A \sharp_{1-\varepsilon} C\right)\right] A^{\frac{1}{2}-\varepsilon} \\
& \leq A^{\frac{1}{2}-\varepsilon}\left[C^{\beta \varepsilon} \sharp_{-\beta}\left(C^{\beta} \sharp_{1-\varepsilon} C\right)\right] A^{\frac{1}{2}-\varepsilon},
\end{aligned}
$$

using the joint monotonicity of matrix geometric means. Since a direct computation yields

$$
C^{\beta \varepsilon} \sharp_{-\beta}\left(C^{\beta} \sharp_{1-\varepsilon} C\right)=C^{\beta(2 \varepsilon-1)}
$$

and by Löwner-Heinz inequality and $0 \leq 1-2 \varepsilon \leq 1, C^{-\alpha} \leq A^{-1}$ implies $C^{-\beta(1-2 \varepsilon)} \leq$ $A^{-(1-2 \varepsilon)}$ and thus we get

$$
A^{r} \natural_{\beta} B^{r} \leq A^{\frac{1}{2}-\varepsilon} C^{\beta(2 \varepsilon-1)} A^{\frac{1}{2}-\varepsilon} \leq A^{\frac{1}{2}-\varepsilon} A^{-1+2 \varepsilon} A^{\frac{1}{2}-\varepsilon}=I .
$$

Therefore (3.4) is proved in the case of $\frac{1}{2} \leq r \leq 1$.

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When $0<r<\frac{1}{2}$, writing $r=2^{-k}(1-\varepsilon)$ with $k \in \mathbb{N}$ and $0 \leq \varepsilon \leq \frac{1}{2}$, and repeating the argument above we have

$$
\begin{aligned}
\lambda_{1}\left(A^{r} \natural_{\beta} B^{r}\right) & \leq \lambda_{1}\left(A^{2^{-(k-1)}(1-\varepsilon)} \natural_{\beta} B^{2^{-(k-1)}(1-\varepsilon)}\right)^{\frac{1}{2}} \\
& \vdots \\
& \leq \lambda_{1}\left(A^{1-\varepsilon} \natural_{\beta} B^{1-\varepsilon}\right)^{2^{-k}} \\
& \leq \lambda_{1}\left(A \natural_{\beta} B\right)^{r}
\end{aligned}
$$

and so the proof is complete.
By Theorem 3.1, we have the following results:
Theorem 3.2. Let $A$ and $B$ be positive definite matrices and $\|\cdot\| \|$ any unitarily invariant norm, and $\beta \in[-1,0)$. If $f$ is a continuous non-decreasing function on $[0, \infty)$ such that $f(0) \geq 0$ and $f\left(e^{t}\right)$ is convex, then

$$
\left\|f\left(A^{r} \natural_{\beta} B^{r}\right)\right\| \leq\left\|f\left(\left(A \natural_{\beta} B\right)^{r}\right)\right\| \quad \text { for all } 0<r \leq 1
$$

In particular,

$$
\left\|A^{r} \natural_{\beta} B^{r}\right\| \leq\left\|\left(A \natural_{\beta} B\right)^{r}\right\| \quad \text { for all } 0<r \leq 1
$$

or equivalently

$$
\begin{gathered}
\left\|\left(A দ_{\beta} B\right)^{r}\right\| \leq\left\|\left(A^{r} দ_{\beta} B^{r}\right)\right\| \quad \text { for all } r \geq 1, \\
\left\|\left(A^{q} দ_{\beta} B^{q}\right)^{\frac{1}{q}}\right\| \leq\left\|\left(A^{p} দ_{\beta} B^{p}\right)^{\frac{1}{p}}\right\| \quad \text { for all } 0<q \leq p .
\end{gathered}
$$

Proof. By [9, Proposition 4.4.13], if $A \prec_{w(\log )} B$ for positive definite matrices $A$ and $B$ and $f$ is a continuous non-decreasing function on $[0, \infty)$ such that $f(0) \geq 0$ and $f\left(e^{t}\right)$ is convex, then $f(A) \prec_{w} f(B)$ and so $\|f(A)\| \leq\|f(B)\|$. Hence Theorem 3.2 follows from Theorem 3.1.

Corollary 3.3. For every positive definite matrices $A, B>0$ and $\beta \in[-1,0)$,

$$
A \natural_{\beta} B \leq I \quad \text { implies } \quad A^{r} দ_{\beta} B^{r} \leq I \quad \text { for all } 0<r \leq 1
$$

4 Norm inequalities for quasi geometric mean. In this section, we show the main norm inequalities related to the quasi geometric mean for positive definite matrices. By [5, Lemma 5.5], we have the following quasi-geometric mean version of the Lie-Trotter formula: If $A$ and $B$ are positive definite matrices, then for each $\beta \in[-1,0)$

$$
\begin{equation*}
A \diamond_{\beta} B=\lim _{p \rightarrow 0}\left(A^{p} \natural_{\beta} B^{p}\right)^{\frac{1}{p}} \tag{4.1}
\end{equation*}
$$

and so for each $\beta \in[-1,0)\left\|\left(A^{p} দ_{\beta} B^{p}\right)^{\frac{1}{p}}\right\| \|$ decreases to $\left\|A \diamond_{\beta} B\right\|$ as $p \downarrow 0$. Hence we have the following norm inequality for the quasi geometric mean of negative power:

Theorem 4.1. Let $A$ and $B$ be positive definite matrices. Then for every unitarily invariant norm

$$
\left\|A \diamond_{\beta} B\right\| \leq\left\|A \natural_{\beta} B\right\| \quad \text { for all } \beta \in[-1,0)
$$

Proof. By Theorem 3.2, it follows that

$$
\left\|\left(A^{q} দ_{\beta} B^{q}\right)^{\frac{1}{q}}\right\| \leq \leq\left\|\left(A^{p} দ_{\beta} B^{p}\right)^{\frac{1}{p}}\right\| \quad \text { for all } 0<q<p
$$

and as $q \rightarrow 0$ and $p=1$ we have the desired inequality by (4.1).

Theorem 4.2. Let $A$ and $B$ be positive definite matrices. Then for every unitarily invariant norm

$$
\begin{equation*}
\left\|A দ_{\beta} B\right\| \leq\left\|A^{1-\beta} B^{\beta}\right\| \| \quad \text { for all } \beta \in\left[-1,-\frac{1}{2}\right] \text {. } \tag{4.2}
\end{equation*}
$$

Proof. For the matrix norm $\|\cdot\|$, by the Araki-Cordes inequality $\left\|B^{t} A^{t} B^{t}\right\| \leq\left\|(B A B)^{t}\right\|$ for all $t \in[0,1]$, we have for $-1 \leq \beta \leq-\frac{1}{2}$

$$
\begin{aligned}
\left\|A দ_{\beta} B\right\| & =\left\|A^{\frac{1}{2}}\left(A^{-\frac{1}{2}} B A^{-\frac{1}{2}}\right)^{\beta} A^{\frac{1}{2}}\right\| \\
& =\left\|A^{\frac{1}{2}}\left(A^{\frac{1}{2}} B^{-1} A^{\frac{1}{2}}\right)^{-\beta} A^{\frac{1}{2}}\right\| \\
& \leq\left\|A^{-\frac{1}{2 \beta}} A^{\frac{1}{2}} B^{-1} A^{\frac{1}{2}} A^{-\frac{1}{2 \beta}}\right\|^{-\beta} \quad \text { by } \frac{1}{2} \leq-\beta \leq 1 \\
& =\left\|A^{\frac{\beta-1}{2 \beta}} B^{-1} A^{\frac{\beta-1}{2 \beta}}\right\|^{-\beta} \\
& \leq\left\|A^{1-\beta} B^{2 \beta} A^{1-\beta}\right\|^{\frac{1}{2}} \quad \text { for } \frac{1}{2} \leq-\frac{1}{2 \beta} \leq 1 \\
& =\left\|\left(A^{1-\beta} B^{2 \beta} A^{1-\beta}\right)^{\frac{1}{2}}\right\|
\end{aligned}
$$

and this implies

$$
\lambda_{1}\left(A \natural_{\beta} B\right) \leq \lambda_{1}\left(\left(A^{1-\beta} B^{2 \beta} A^{1-\beta}\right)^{\frac{1}{2}}\right)=\lambda_{1}\left(\left|B^{\beta} A^{1-\beta}\right|\right) .
$$

Replacing $A$ and $B$ by (5) of Lemma 2.1, we obtain

$$
\prod_{i=1}^{k} \lambda_{i}\left(A \natural_{\beta} B\right) \leq \prod_{i=1}^{k} \lambda_{i}\left(\left|B^{\beta} A^{1-\beta}\right|\right) \quad \text { for } k=1, \ldots, n
$$

Hence we have the weak $\log$ majorization $A \natural_{\beta} B \prec_{w(\log )}\left|B^{\beta} A^{1-\beta}\right|$ and this implies

$$
\left\|A \natural_{\beta} B\right\| \leq\| \|\left|B^{\beta} A^{1-\beta}\right|\| \|=\left\|B^{\beta} A^{1-\beta}\right\|\|=\| A^{1-\beta} B^{\beta}\| \|
$$

for every unitarily invariant norm and so we have the desired inequality (4.2).
Remark 4.3. In Theorem 4.2, the inequality $\left\|A \natural_{\beta} B\right\| \leq\left\|A^{1-\beta} B^{\beta}\right\|$ does not always hold for $-1 / 2<\beta<0$. In fact, if we put $\beta=-\frac{1}{3}, A=\left(\begin{array}{ll}2 & 1 \\ 1 & 2\end{array}\right)$ and $B=\left(\begin{array}{ll}2 & 1 \\ 1 & 1\end{array}\right)$, then we have the matrix norm $\left\|A \natural_{-\frac{1}{3}} B\right\|=3.385$ and $\left\|A^{\frac{4}{3}} B^{-\frac{1}{3}}\right\|=3.375$, and so $\left\|A \natural_{\beta} B\right\|>\left\|A^{1-\beta} B^{\beta}\right\|$.
Theorem 4.4. Let $A$ and $B$ be positive definite matrices. Then for every unitarily invariant norm

$$
\begin{equation*}
\left\|A \diamond_{\beta} B\right\| \leq\left\|\left(B^{\frac{\beta q}{2}} A^{(1-\beta) q} B^{\frac{\beta q}{2}}\right)^{\frac{1}{q}}\right\| \leq\left\|A \natural_{\beta} B\right\| \quad \text { for } 0<q \leq \frac{1}{2} \text { and } \beta \in[-1,0) \tag{4.3}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\|A দ_{\beta} B\right\| \leq\left\|A^{1-\beta} B^{\beta}\right\|\|\leq\|\left(B^{\frac{\beta p}{2}} A^{(1-\beta) p} B^{\frac{\beta p}{2}}\right)^{\frac{1}{p}} \| \quad \text { for } p \geq 2 \text { and } \beta \in\left[-1,-\frac{1}{2}\right] \tag{4.4}
\end{equation*}
$$

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Proof. Since the first inequality in (4.3) follows from the Lie-Trotter formula (4.1), we show the second inequality in (4.3). By Theorem 3.2, we have $\left\|\left(A^{r} \natural_{\beta} B^{r}\right)^{\frac{1}{r}}\right\| \leq\left\|A \natural_{\beta} B\right\|$ for all $0<r \leq 1$. For the matrix norm $\|\cdot\|$, by the Araki-Cordes inequality $\left\|(B A B)^{t}\right\| \leq\left\|B^{t} A^{t} B^{t}\right\|$ for all $t \geq 1$, we have for $0<r \leq 1$

$$
\begin{aligned}
\left\|A \mathfrak{\natural}_{\beta} B\right\| & \geq\left\|A^{r} \mathfrak{h}_{\beta} B^{r}\right\|^{\frac{1}{r}} \\
& =\left\|B^{r} \mathfrak{\natural}_{1-\beta} A^{r}\right\|^{\frac{1}{r}} \\
& =\left\|B^{\frac{r}{2}}\left(B^{-\frac{r}{2}} A^{r} B^{-\frac{r}{2}}\right)^{1-\beta} B^{\frac{r}{2}}\right\|^{\frac{1}{r}} \\
& \geq\left\|B^{\frac{\beta r}{2(1-\beta)}} A^{r} B^{\frac{\beta r}{2(1-\beta)}}\right\|^{\frac{1-\beta}{r}} \quad \text { by } 0<\frac{1}{1-\beta}<1 \\
& \geq\left\|B^{\frac{\beta r}{4}} A^{\frac{(1-\beta) r}{2}} B^{\frac{\beta r}{4}}\right\|^{\frac{2}{r}} \quad \text { by } \frac{1}{2}<\frac{1-\beta}{2} \leq 1 .
\end{aligned}
$$

If we put $q=\frac{r}{2}$, then we have $\left\|\left(B^{\frac{\beta q}{2}} A^{(1-\beta) q} B^{\frac{\beta q}{2}}\right)^{\frac{1}{q}}\right\| \leq \| A$ म $_{\beta} B \|$ for $0<q \leq \frac{1}{2}$ and this implies

$$
\lambda_{1}\left(\left(B^{\frac{\beta q}{2}} A^{(1-\beta) q} B^{\frac{\beta q}{2}}\right)^{\frac{1}{q}}\right) \leq \lambda_{1}\left(A \mathfrak{h}_{\beta} B\right) .
$$

Replacing $A$ and $B$ by (5) of Lemma 2.1, we obtain

$$
\prod_{i=1}^{k} \lambda_{1}\left(\left(B^{\frac{\beta q}{2}} A^{(1-\beta) q} B^{\frac{\beta q}{2}}\right)^{\frac{1}{q}}\right) \leq \prod_{i=1}^{k} \lambda_{1}\left(A \natural_{\beta} B\right) \quad \text { for } k=1, \ldots, n,
$$

which gives the second inequality in (4.3).
Next, for $s \geq 1$, it follows from Theorem 3.2 that $\left\|A \natural_{\beta} B\right\| \leq\left\|\left(A^{s} দ_{\beta} B^{s}\right)^{\frac{1}{s}}\right\|$. For the matrix norm $\|\cdot\|$, we have

$$
\begin{aligned}
\left\|A \mathfrak{h}_{\beta} B\right\| & \leq\left\|A^{s} দ_{\beta} B^{s}\right\|^{\frac{1}{s}} \\
& =\left\|A^{\frac{s}{2}}\left(A^{\frac{s}{2}} B^{-s} A^{\frac{s}{2}}\right)^{-\beta} A^{\frac{s}{2}}\right\|^{\frac{1}{s}} \\
& \leq\left\|A^{\frac{-(1-\beta) s}{2 \beta}} B^{-s} A^{\frac{-(1-\beta) s}{2 \beta}}\right\|^{-\frac{\beta}{s}} \quad \text { by } \frac{1}{2} \leq-\beta \leq 1 \\
& \leq\left\|A^{(1-\beta) s} B^{2 \beta s} A^{(1-\beta) s}\right\|^{\frac{1}{2 s}} \quad \text { by } \frac{1}{2} \leq-\frac{1}{2 \beta} \leq 1 .
\end{aligned}
$$

If we put $p=2 s$, then we have

$$
\begin{aligned}
\| A \text { n }_{\beta} B \| & \leq\left\|A^{\frac{(1-\beta) p}{2}} B^{\beta p} A^{\frac{(1-\beta) p}{2}}\right\|^{\frac{1}{p}} \\
& =\operatorname{spr}\left(A^{\frac{(1-\beta) p}{2}} B^{\beta p} A^{\frac{(1-\beta) p}{2}}\right)^{\frac{1}{p}} \\
& =\operatorname{spr}\left(B^{\frac{\beta p}{2}} A^{(1-\beta) p} B^{\frac{\beta p}{2}}\right)^{\frac{1}{p}} \\
& =\left\|B^{\frac{\beta p}{2}} A^{(1-\beta) p} B^{\frac{\beta p}{2}}\right\|^{\frac{1}{p}} \\
& =\left\|\left(B^{\frac{\beta p}{2}} A^{(1-\beta) p} B^{\frac{\beta p}{2}}\right)^{\frac{1}{p}}\right\|
\end{aligned}
$$

for $p \geq 2$, where $\operatorname{spr}(X)$ is the spectral radius of $X$. By the argument similar to above, we have the inequality (4.4).

Let $A$ and $B$ be positive definite matrices in $\mathbb{M}_{n}$ and $\beta \in[-1,0)$. Since there is the case that $(1-\beta) A+\beta B$ is not positive semidefinite, we have no relation between $\|(1-\beta) A+\beta B\|$ and $\left\|A \natural_{\beta} B\right\|$ though $(1-\beta) A+\beta B \leq A দ_{\beta} B$. Suppose that $A \geq B$. Then $(1-\beta) A^{p}+$ $\beta B^{p}$ is positive definite for all $p \in(0,1]$. In particular $0<(1-\beta) A+\beta B \leq A \natural_{\beta} B$ and so $\|(1-\beta) A+\beta B\| \leq\left\|A \natural_{\beta} B\right\|$ for every unitarily invariant norm. Thus under the assumption $A \geq B$, we consider the refinement of this norm inequality. For this, we need the following result due to J. I. Fujii [6]: A real valued continuous function $f$ on an interval $J$ is matrix concave if and only if

$$
\begin{equation*}
f((1-\beta) H+\beta K) \leq(1-\beta) f(H)+\beta f(K) \tag{4.5}
\end{equation*}
$$

for all Hermitian matrices $H$ and $K$ with $\sigma(H), \sigma(K)$ and $\sigma((1-\beta) H+\beta K) \subset J$ and $\beta \in[-1,0)$.

Let $0<q<p \leq 1$. Then the function $f(t)=t^{\frac{q}{p}}$ on $[0, \infty)$ is matrix concave and by (4.5)

$$
\begin{equation*}
\left((1-\beta) A^{p}+\beta B^{p}\right)^{\frac{q}{p}} \leq(1-\beta) A^{q}+\beta B^{q} \tag{4.6}
\end{equation*}
$$

Note that $(1-\beta) A^{p}+\beta B^{p}>0$ for all $p \in(0,1]$ since $A \geq B$. This implies that

$$
\lambda_{i}\left((1-\beta) A^{p}+\beta B^{p}\right)^{\frac{q}{p}} \leq \lambda_{i}\left((1-\beta) A^{q}+\beta B^{q}\right) \quad \text { for all } i=1, \ldots, n
$$

Taking $q$-th roots of both sides, we obtain

$$
\lambda_{i}\left((1-\beta) A^{p}+\beta B^{p}\right)^{\frac{1}{p}} \leq \lambda_{i}\left((1-\beta) A^{q}+\beta B^{q}\right)^{\frac{1}{q}} \quad \text { for all } i=1, \ldots, n
$$

and so $\left\|\left\|\left((1-\beta) A^{p}+\beta B^{p}\right)^{\frac{1}{p}}\right\|\right\|$ is a decreasing function of $p$.
On the other hand, taking the logarithm of both sides in (4.6) and by (4.5), we obtain

$$
\begin{aligned}
\log \left((1-\beta) A^{p}+\beta B^{p}\right)^{\frac{1}{p}} & \leq \frac{1}{q} \log \left((1-\beta) A^{q}+\beta B^{q}\right) \\
& \leq(1-\beta) \log A+\beta \log B
\end{aligned}
$$

and this implies

$$
\lambda_{i}\left(\log \left((1-\beta) A^{p}+\beta B^{p}\right)^{\frac{1}{p}}\right) \leq \lambda_{i}((1-\beta) \log A+\beta \log B) \quad \text { for all } i=1, \ldots, n
$$

Taking the exponent of both sides, we obtain

$$
\lambda_{i}\left((1-\beta) A^{p}+\beta B^{p}\right)^{\frac{1}{p}} \leq \lambda_{i}\left(e^{(1-\beta) \log A+\beta \log B}\right) \quad \text { for all } i=1, \ldots, n
$$

and so

$$
\left\|\left((1-\beta) A^{p}+\beta B^{p}\right)^{\frac{1}{p}}\right\| \leq\left\|e^{(1-\beta) \log A+\beta \log B}\right\|
$$

for all $p \in(0,1]$. Summing up, we obtain the following result:
Theorem 4.5. Let $A$ and $B$ be positive definite matrices in $\mathbb{M}_{n}$ such that $A \geq B$ and $\beta \in[-1,0)$. Then for every unitarily invariant norm

$$
\|(1-\beta) A+\beta B\| \leq\| \|\left((1-\beta) A^{p}+\beta B^{p}\right)^{\frac{1}{p}}\|\leq\| A \diamond_{\beta} B\|\leq\| A দ_{\beta} B \|
$$

for all $p \in(0,1]$.

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Finally, as an application, we show a refinement of the generalized Golden-Thompson inequality in terms of the quasi geometric means. Let $H$ and $K$ be Hermitian matrices. The Golden-Thompson trace inequality is

$$
\operatorname{Tr}\left[e^{H+K}\right] \leq \operatorname{Tr}\left[e^{H} e^{K}\right]
$$

Hiai-Petz [10] proved the complemented Golden-Thompson inequality:

$$
\left\|e^{H} \sharp_{\alpha} e^{K}\right\| \leq \leq\left\|e^{(1-\alpha) H+\alpha K}\right\| \| \quad \text { for all } \alpha \in[0,1]
$$

for every unitarily invariant norm. By Theorem 4.1 and Theorem 4.2, we have a refinement of the Golden-Thompson inequality in terms of the quasi geometric means:

$$
\left\|e^{(1-\beta) H+\beta K}\right\| \leq\left\|e^{H} \natural_{\beta} e^{K}\right\| \leq\left\|e^{(1-\beta) H} e^{\beta K}\right\|
$$

for all $\beta \in\left[-1,-\frac{1}{2}\right]$ and so

$$
\operatorname{Tr}\left[e^{H+K}\right] \leq \operatorname{Tr}\left[e^{\frac{1}{1-\beta} H} দ_{\beta} e^{\frac{1}{\beta} K}\right] \leq \operatorname{Tr}\left[e^{H} e^{K}\right]
$$

In particular, if we put $\alpha=\frac{1}{2}$ and $\beta=-\frac{1}{2}$, then we have

$$
\operatorname{Tr}\left[e^{2 H} \sharp_{\frac{1}{2}} e^{2 K}\right] \leq \operatorname{Tr}\left[e^{H+K}\right] \leq \operatorname{Tr}\left[e^{\frac{2}{3} H} \natural_{-\frac{1}{2}} e^{-2 K}\right] \leq \operatorname{Tr}\left[e^{H} e^{K}\right] .
$$

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# ON THE BANG-BANG PRINCIPLE FOR DIFFERENTIAL INCLUSIONS IN A REFLEXIVE SEPARABLE BNANACH SPACE 

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#### Abstract

In this paper, we consider the relation existing between the solutions of the following differential inclusions: (I) $\dot{x} \in \Gamma(t, x), x(0)=0$ and (II) $\dot{x} \in \operatorname{ext} \Gamma(t, x), x(0)=0$ defined on a reflexive separable Banach space. In particular, we establish the sufficient conditions which guarantee the set of solutions of (II) is dense in the set of solutions of (I) with respect to the (weak) uniformly continuous topology.


Let $(\mathfrak{X},\|\cdot\|)$ be a real reflexive separable Banach space and $T$ be a positive real number. Let $\Gamma:[0, T] \times \mathfrak{X} \rightarrow \mathfrak{X}$ be a correspondence (=multi-valued function). We consider a relation existing between the sets of solutions of the following differential inclusions:
(I) $\dot{x} \in \Gamma(t, x), x(0)=0$, and
(II) $\dot{x} \in \operatorname{ext} \Gamma(t, x), x(0)=0$,
where $\operatorname{ext} A$ stands for the weak-closure of the extreme points of $A$. By a solution of (I) and (II), we mean an absolutely continuous function $x:[0, T] \rightarrow \mathfrak{X}$ that satisfies $\dot{x} \in \Gamma(t, x(t))$ a.e. in $t \in[0, T]$ and $x(0)=0$ in the case of (I) and $\dot{x} \in \operatorname{ext} \Gamma(t, x(t))$ a.e. in $t \in[0, T]$ and $x(0)=0$ in the case of (II). We denote by $\mathcal{R}$ and $\mathcal{R}_{*}$ the set solutions of (I) and (II) respectively. Tateishi $[5,6]$ established the existence of solutions of the differential inclusions (I) under the following assumptions:
(i) $\Gamma$ is nonempty and weakly compact-valued, i.e., $\Gamma(t, x)$ is nonempty and weakly compact for each $(t, x) \in[0, T] \times \mathfrak{X}$,
(ii) for each fixed $t \in[0, T]$, the correspondence $t \rightarrow \Gamma(t, x)$ is continuous with respect to the weak topology for $\mathfrak{X}$,
(iii) for each fixed $x \in \mathfrak{X}$, the correspondence $t \rightarrow \Gamma(t, x)$ is measurable, and
(iv) there exists $M>0$ such that $\sup \{\|y\| \mid y \in \Gamma(t, x), t \in[0, T], x \in \mathfrak{X}\} \leq M$.

Furthermore, Tateishi $[6,7]$ examined the relations existing between the solutions set of (I) and (III): $\dot{x} \in \overline{\mathrm{co}} \Gamma(t, x), x(0)=0$. The aim of this paper is to establish the relation between the sets of solutions (I) and (II). Bressan [1, 2] established the existence of solutions of both of the problems and obtained the closure result $\mathcal{R}=\overline{\mathcal{R}}_{*}$ in the case that $\mathfrak{X}$ is a finite dimensional space. In this paper, we generalize his theorem to infinite dimensional spaces.

[^1]
## 1. Preliminaries

In this section, we offer some notations and lemmata used in this paper.
Let $(\mathfrak{X},\|\cdot\|)$ be a reflexive separable Banach space with $\mathfrak{X}^{*}$ its dual. We denote by $\mathfrak{X}^{w}$ the space $\mathfrak{X}$ endowed with the weak topology. Let $\mathfrak{S}=\{x \in \mathfrak{X} \mid\|x\| \leq \max (M T, M)\}$ where $M$ and $T$ are constants which appear in the Introduction section. The set $\mathfrak{S}$ endowed with the relative topology of $\mathfrak{X}^{w}$ is denoted by $\mathfrak{S}^{w}$. The following proposition is from Larman and Rogers [4, Theorem 2].
Proposition 1. Let E be a Hausdorff locally convex topological vector space. Let $X$ be a metrizable compact subset of $E$. Let $V$ be the linear subspace generated by the set $\overline{\text { co }} X$. Then it is possible to introduce a norm on $V$ so that the relative topologies of $\overline{\operatorname{co}} X$, as a subset of $E$, and as a subset of the normed space $V$, coincide.
$\mathfrak{X}^{w}$ is a Hausdorff locally convex topological vector space and $\mathfrak{S}^{w}$ is a metrizable and compact subset of $\mathfrak{X}^{w}$. Furthermore, the linear subspace generated by $\mathfrak{S}^{w}$ is the whole space $\mathfrak{X}$. Hence, we can, by the above proposition, introduce a norm $\|\cdot\|^{w}$ on $\mathfrak{X}^{w}$ so that the topology on $\mathfrak{S}^{w}$ and the relative topology as a subset of the normed vector space $\left(\mathfrak{X},\|\cdot\|^{w}\right)$ coincide.

We denote by $h$ the Hausdorff distance on $\mathfrak{S}^{w}$ induced by $\|\cdot\|^{w}$, that is, $h\left(A, A^{\prime}\right)=$ $\max \left\{\sup _{x \in A^{\prime}} d(x, A), \sup _{x \in A} d\left(x, A^{\prime}\right)\right\}$ for any closed subsets $A, A^{\prime}$ of $\mathfrak{S}^{w}$, where $d(x, A)=$ $\inf \left\{\|x-y\|^{w} \mid y \in A\right\}$. For $A \subset \mathfrak{S}$ and $\alpha>0$, we set $B[A, \alpha]=\{x \in \mathfrak{S} \mid d(x, A)<\alpha\}$. We denote by $\mu$, the Lebesgue measure defined on the interval $[0, T]$.
Lemma 1. Let $(\mathfrak{X},\|\cdot\|)$ be a normed linear space and $\Gamma: \mathfrak{X} \rightarrow \mathfrak{X}$ be convex, compact-valued and continuous. Then the map $\operatorname{ext} \Gamma: \mathfrak{X} \rightarrow \mathfrak{X}$ is lower hemi-continuous.

Proof. Let $x_{0}, y_{0} \in \mathfrak{X}$ with $y_{0} \in \operatorname{ext} \Gamma\left(x_{0}\right)$ and $\left\{x_{n}\right\}$ be a sequence which converges to $x_{0}$. We must show that for some subsequence $x_{n^{\prime}}$ of $x_{n}$ and some $y_{n^{\prime}} \in \operatorname{ext} \Gamma\left(x_{n^{\prime}}\right)$, we have $y_{n^{\prime}} \rightarrow y_{0}$. Since $\Gamma$ is continuous, there exists a sequence $y_{n} \in \Gamma\left(x_{n}\right)$ such that $y_{n} \rightarrow y_{0}$. Since $\Gamma$ is compact and convex-valued, the Krein-Milman theorem implies that $\Gamma\left(x_{n}\right)=\overline{\operatorname{co}} \operatorname{ext} \Gamma\left(x_{n}\right)$. Hence, for each $n \in \mathbb{N}$, there exists $\alpha_{n}^{i} \geq 0, \sum_{i} \alpha_{n}^{i}=1(i \in \mathbb{N})$, where only finitely many $\alpha_{n}^{i}$ are not equal to zero, and $z_{n}^{i} \in \operatorname{ext} \Gamma\left(x_{n}\right)$ such that $\left\|y_{n}-\sum_{i} \alpha_{n}^{i} z_{n}^{i}\right\| \leq 1 / n$. Let $y_{n}^{i} \in \Gamma\left(x_{0}\right)$ be such that $\left\|z_{n}^{i}-y_{n}^{i}\right\| \leq h\left(\Gamma\left(x_{n}\right), \Gamma\left(x_{0}\right)\right)$, where $h$ is the Hausdorff metric defined by $\|\cdot\|$. Since $\Gamma\left(x_{0}\right)$ is compact, there exist, for each fixed $i$, converging subsequences $y_{n^{\prime}}^{i}$ to $y_{0}^{i}$ and $\alpha_{n^{\prime}}^{i}$ to $\alpha_{0}^{i}$. Then $\sum_{i} \alpha_{0}^{i} y_{0}^{i}=y_{0}$ and since $y_{0}$ is an extreme point of $\Gamma\left(x_{0}\right)$, we have each $y_{0}^{i}$ is equal to $y_{0}$ for all $i$ with $\alpha_{0}^{i}>0$. Let $i^{*}$ be such that $\alpha_{0}^{i^{*}}>0$. Then $\lim \sup _{n^{\prime}}\left\|z_{n^{\prime}}^{i^{*}}-y_{0}\right\| \leq \lim \sup _{n^{\prime}}\left\|z_{n^{\prime}}^{i^{*}}-y_{n^{\prime}}^{i^{*}}\right\|+\lim \sup _{n^{\prime}}\left\|y_{n^{\prime}}^{i^{*}}-y_{0}\right\|=0$. Hence $y_{n^{\prime}}$ also converges to $y_{0}$ and this completes the proof.

Lemma 2. Let $F:[0, T] \times \mathfrak{X}^{w} \rightarrow \mathfrak{X}^{w}$ be lower hemi-continuous and $V \subset \mathfrak{X}^{w}$ be open. Then the correspondence $H:[0, T] \times \mathfrak{X}^{w} \rightarrow \mathfrak{X}^{w}$ defined by $H(t, x)=\overline{F(t, x) \cap V}$ is lower hemi-continuous, where $\bar{A}$ stands for the closure with respect to the weak topology of $\mathfrak{X}$.
Proof. Let $K$ be a weakly closed subset of $\mathfrak{X}$. Then we have the following implications:

$$
H(t, x) \subset K \Leftrightarrow F(t, x) \cap V \subset K \Leftrightarrow F(t, x) \subset K \cup V^{c} .
$$

Since $K \cup V^{c}$ is weakly closed and $F$ is lower hemi-continuous, the set $\{(t, x) \mid H(t, x) \subset$ $K\}=\left\{(t, x) \mid F(t, x) \subset K \cup V^{c}\right\}$ is closed in $[0, T] \times \mathfrak{X}^{w}$. It follows that $H$ is lower hemi-continuous.

Lemma 3. Let $I_{0} \subset[0, T]$ be a measurable set with $\mu\left(I_{0}\right)=\sigma$ and let $M$ and $\epsilon$ be given positive real numbers. Then the solution $\psi:[0, T] \rightarrow \mathbb{R}$ of the differential equation

$$
\begin{equation*}
\dot{\psi}(t)=\psi(t)+2 M \chi_{I_{0}}(t)+4 \epsilon, \psi(0)=0 \tag{1}
\end{equation*}
$$

is positive, monotonically increasing and satisfies the following inequality:

$$
\psi(T) \leq 2 M \sigma e^{T}+4 \epsilon\left(e^{T}-1\right)
$$

Proof. It is easy to verify that the solution $\psi$ is positive and monotonically increasing. By calculating the solution $\psi$ of (1) directly, we obtain

$$
\psi(T)=2 M \cdot \int_{0}^{T} \chi_{I_{0}}(s) e^{(T-s)} d s+4 \epsilon \int_{0}^{T} e^{(T-s)} d s \leq 2 M \sigma e^{T}+4 \epsilon\left(e^{T}-1\right) .
$$

## 2. Main Theorem

Theorem 1. Let $\Gamma:[0, T] \times \mathfrak{X} \rightarrow \mathfrak{X}$ be a correspondence which satisfies the conditions:
(i) $\Gamma$ is convex and weakly compact-valued, that is $\Gamma(t, x)$ is convex and weakly compact for each $(t, x) \in[0, T] \times \mathfrak{X}$,
(ii) $\Gamma$ is continuous, where $\mathfrak{X}$ is endowed with the weak topology.
(iii) $h(\Gamma(t, x), \Gamma(t, y)) \leq\|x-y\|^{w}$, and
(iv) there exists $M>0$ such that $\sup \{\|y\| \mid y \in \Gamma(t, x), t \in[0, T], x \in \mathfrak{X}\} \leq M$.

Then $\mathcal{R}=\overline{\mathcal{R}_{*}}$, that is the set of solutions of (II) is dense in the set of solutions of (I) with respect to the (weak) uniform convergence topology.

Proof. Step 1. Let $v$ be a solution of (I) and let $\epsilon$ be a positive real number. Then there exists, an open subset $I_{0}$ of $[0, T]$ with $\mu\left(I_{0}\right)<\epsilon$ such that, for all $t \in I_{1}=[0, T] \cap I_{0}^{c}, \dot{v}(t)$ exists and lies in $\Gamma(t, v(t))$, and the restriction $\left.\dot{v}\right|_{I_{1}}$ of $\dot{v}$ to $I_{1}$ is continuous. We may also assume that $\left[0, \tau_{0}\right] \subset I_{0}$ for some $\tau_{0}>0$.

Step 2. Let $\mathfrak{M}$ be the set of $\{u, \tau\}$ of an absolutely continuous mapping $u$ and a positive constant $0 \leq \tau \leq T$ such that $u$ is defined on the closed interval $[0, \tau]$ and satisfies $\dot{u} \in \operatorname{ext} \Gamma(t, u(t))$ a.e. in $t \in[0, \tau], u(0)=0$,

$$
\begin{equation*}
\|u(\tau)-v(\tau)\|^{w} \leq \psi(\tau), \text { and } \tag{2}
\end{equation*}
$$

$$
\begin{equation*}
\|u(t)-v(t)\|^{w} \leq \psi(t)+2 M \epsilon \text { for all } t \in[0, \tau], \tag{3}
\end{equation*}
$$

where $\psi$ is a solution of (1).
Step 3. Since $\left[0, \tau_{0}\right] \subset I_{0}$, the pair $\left\{u, \tau_{0}\right\}$ for every solution $u$ of (II) satisfies the above properties, thus the set $\mathfrak{M}$ is nonempty. Let us define a partial ordering $\precsim \mathfrak{M}$ on $\mathfrak{M}$ by $\left(u_{1}, \tau_{1}\right) \precsim \mathfrak{M}\left(u_{2}, \tau_{2}\right) \Leftrightarrow \tau_{1} \leq \tau_{2}$ and $u_{2}$ is an extension of $u_{1}$. Then Zorn's lemma implies that there exists a maximal element $\left(u^{*}, \tau^{*}\right)$ of $\mathfrak{M}$.

Step 4. Since $\epsilon>0$ is arbitrary, the equations (2) and (3) imply that the solution $u^{*}$ of (II) can be arbitrarily near to the solution $v$ with respect to the (weak) uniform convergence topology on $\left[0, \tau^{*}\right]$. In the following two steps, we show that $\tau^{*}$ obtained in Step 3 equals $T$. In this step, we consider the case $\tau^{*} \in I_{0}$. Then, since $I_{0}$ is open, there exists a positive number $\delta$ such that $\left[\tau^{*}, \tau^{*}+\delta\right] \subset I_{0}$. Then, we have an absolutely continuous function $w:\left[\tau^{*}, \tau^{*}+\delta\right] \rightarrow X$ satisfying $\dot{w}(t) \in \operatorname{ext} \Gamma(t, w(t))$ for $t \in\left[\tau^{*}, \tau^{*}+\delta\right], w\left(\tau^{*}\right)=u^{*}\left(\tau^{*}\right)$. Let us define $w^{*}:\left[0, \tau^{*}+\delta\right] \rightarrow \mathfrak{X}$ by

$$
w^{*}(t)= \begin{cases}u^{*}(t) & \text { for } t \in\left[0, \tau^{*}\right] \\ w(t) & \text { for } t \in\left[\tau^{*}, \tau^{*}+\delta\right]\end{cases}
$$

Then, for $t \in\left[\tau^{*}, \tau^{*}+\delta\right]$, we obtain the estimation:

$$
\begin{aligned}
& \left\|w^{*}(t)-v(t)\right\|^{w} \\
& \leq\left\|w^{*}\left(\tau^{*}\right)-v\left(\tau^{*}\right)\right\|^{w}+\int_{\tau^{*}}^{t}\left\|\dot{w}^{*}(s)-\dot{v}(s)\right\|^{w} d s \\
& \leq\left\|w^{*}\left(\tau^{*}\right)-v\left(\tau^{*}\right)\right\|^{w}+\int_{\tau^{*}}^{t} 2 M \chi_{I_{0}}(s) d s \\
& \leq \psi\left(\tau^{*}\right)+\int_{\tau^{*}}^{t} \dot{\psi}(s) d s=\psi(t)
\end{aligned}
$$

where the first inequality is an immediate consequence of the fundamental theorem of calculus, and the second follows from assumption (iv) and the third from (2) and the definition of $\psi$. Thus $w^{*}$ belongs to $\mathfrak{M}$, which contradicts the maximality of $u^{*}$.

Step 5. In this step, we consider the case $\tau^{*} \in I_{1}$. By assumption (iii) and (2), we have $y^{*} \in \Gamma\left(\tau^{*}, u^{*}\left(\tau^{*}\right)\right)$ such that $\left\|y^{*}-\dot{v}\left(\tau^{*}\right)\right\|^{w} \leq \psi\left(\tau^{*}\right)$. Since $\left.\dot{v}^{*}\right|_{I_{1}}$ is continuous by assumption and $\operatorname{ext} \Gamma: I_{1} \times \mathfrak{X}^{w} \rightarrow \mathfrak{X}^{w}$ is lower hemi-continuous by Lemma 1 , we have a positive constant $0<\delta \leq \epsilon$ such that $\left\|\dot{v}^{*}(t)-\dot{v}^{*}(\tau)\right\|^{w}<\epsilon$ for $t \in I_{1} \cap\left[\tau^{*}, \tau^{*}+\delta\right]$, and

$$
\begin{equation*}
\operatorname{ext} \Gamma\left(\tau^{*}, u^{*}\left(\tau^{*}\right)\right) \subset B[\operatorname{ext} \Gamma(t, x), \epsilon] \text { for } t \in\left[\tau^{*}, \tau^{*}+\delta\right] \text { and } x \in B\left[u^{*}\left(\tau^{*}\right), M \delta\right] \tag{4}
\end{equation*}
$$

By the Krein-Milman theorem, we have: $y^{*} \in \Gamma\left(\tau^{*}, u^{*}\left(\tau^{*}\right)\right)=\overline{\operatorname{co}} \operatorname{ext} \Gamma\left(\tau^{*}, u^{*}\left(\tau^{*}\right)\right)$, where $\overline{\mathrm{co}}$ stands for the closed convex hull of $A$. Thus we obtain, for any $\epsilon>0$, finite points $y_{1}, y_{2}, \ldots, y_{m}$ in $\operatorname{ext} \Gamma\left(\tau^{*}, u^{*}\left(\tau^{*}\right)\right)$ and nonnegative real numbers $\lambda_{1}, \lambda_{2}, \ldots, \lambda_{m}$ with $\sum_{i=1}^{m} \lambda_{i}=1$ such that $\left\|y^{*}-\sum_{i} \lambda_{i} y_{i}\right\|^{w}<\epsilon$. Then we obtain, by Lyapunov's convexity theorem, a set of $m$ measurable partition $J_{0}, J_{1}, \ldots, J_{m}$ such that $t<s$ for $t \in J_{i}, s \in J_{j}$ with $i<j, \cup_{i} J_{i}=I_{1} \cap\left[\tau^{*}, \tau^{*}+\delta\right]$, and $\mu\left(J_{i}\right)=\lambda_{i} \mu\left(I_{1} \cap\left[\tau^{*}, \tau^{*}+\delta\right]\right)$. For $t \in\left[\tau^{*}, \tau^{*}+\delta\right]$, we set

$$
H(t, x)= \begin{cases}\operatorname{ext} \Gamma(t, x) & \text { if } t \in I_{0} \\ \overline{\operatorname{ext} \Gamma(t, x) \cap B\left[y_{i}, \epsilon\right]} & \text { if } t \in J_{i} .\end{cases}
$$

Then by (4), $y_{i} \in \operatorname{ext} \Gamma\left(\tau^{*}, u^{*}\left(\tau^{*}\right)\right) \subset B[\operatorname{ext} \Gamma(t, x), \epsilon]$ for $i=1,2, \ldots, m, t \in\left[\tau^{*}, \tau^{*}+\delta\right]$, and $x \in B\left[u^{*}\left(\tau^{*}\right), M \delta\right]$. It follows that $\operatorname{ext} \Gamma(t, x) \cap B\left[y_{i}, \epsilon\right] \neq \emptyset$ and hence, $H(t, x) \neq \emptyset$ for such pair $(t, x)$. In view of Lemma 2, the restriction of $H$ to each of the product spaces $I_{0} \times B\left[u^{*}(\tau), M \delta\right]$ and $J_{i} \times B\left[u^{*}(\tau), M \delta\right]$ is lower hemi-continuous. Thus $H$ is almost lower hemi-continuous and we have an absolutely continuous function $u_{\delta}:\left[\tau^{*}, \tau^{*}+\delta\right] \rightarrow$
$\mathfrak{X}$ satisfying the following conditions: $u_{\delta}\left(\tau^{*}\right)=u^{*}\left(\tau^{*}\right)$, and $\dot{u}_{\delta}(t) \in H\left(t, u_{\delta}(t)\right)$ a.e. in $\left[\tau^{*}, \tau^{*}+\delta\right]$ (see, e.g., Deimling [3,Theorem 9.3]). Let us define $w^{*}:\left[0, \tau^{*}+\delta\right] \rightarrow \mathfrak{X}$ by

$$
w^{*}(t)= \begin{cases}u^{*}(t) & \text { for } t \in\left[0, \tau^{*}\right] \\ u_{\delta}(t) & \text { for } t \in\left[\tau^{*}, \tau^{*}+\delta\right]\end{cases}
$$

Then, $w^{*}$ can be seen to be an element of $\mathfrak{M}$ as follows. First, we verify that $w^{*}$ satisfies the condition (2) for $\tau=\tau^{*}+\delta$. Setting $I_{0}^{*}=I_{0} \cap\left[\tau^{*}, \tau^{*}+\delta\right], I_{1}^{*}=I_{1} \cap\left[\tau^{*}, \tau^{*}+\delta\right]$, we have

$$
\begin{aligned}
& \left\|w^{*}\left(\tau^{*}+\delta\right)-v\left(\tau^{*}+\delta\right)\right\|^{w} \\
& \leq\left\|w^{*}\left(\tau^{*}\right)-v\left(\tau^{*}\right)\right\|^{w}+\left\|\int_{\tau^{*}}^{\tau^{*}+\delta}\left[\dot{u}_{\delta}(t)-\dot{v}(t)\right] d t\right\|^{w} \\
& \leq \psi\left(\tau^{*}\right)+\int_{I_{0}^{*}}\left\|\dot{u}_{\delta}(t)-\dot{v}(t)\right\|^{w} d t+\left\|\sum_{i=1}^{m} \int_{J_{i}} \dot{u}_{\delta}(t) d t-\int_{I_{1}^{*}} \dot{v}(t) d t\right\|^{w} \\
& \leq \psi\left(\tau^{*}\right)+2 M \mu\left(I_{0}^{*}\right)+\sum_{i=1}^{m} \int_{J_{i}}\left\|\dot{u}_{\delta}(t)-y_{i}\right\|^{w} d t+\mu\left(I_{1}^{*}\right)\left\|\sum_{i=1}^{m} \lambda_{i} y_{i}-y^{*}\right\|^{w} \\
& +\delta\left\|y^{*}-\dot{v}\left(\tau^{*}\right)\right\|^{w}+\epsilon \mu\left(I_{1}^{*}\right) \\
& \leq \psi\left(\tau^{*}\right)+2 M \mu\left(I_{0}^{*}\right)+4 \epsilon \mu\left(I_{1}^{*}\right)+\delta\left\|y^{*}-\dot{v}\left(\tau^{*}\right)\right\|^{w} \\
& \leq \psi\left(\tau^{*}\right)+\int_{\tau^{*}}^{\tau^{*}+\delta}\left(2 M \chi_{I_{0}}(t)+4 \epsilon t\right) d t+\delta \psi\left(\tau^{*}\right) \leq \psi\left(\tau^{*}+\delta\right),
\end{aligned}
$$

where the first inequality is a consequence of the fundamental theorem of calculus. In the second inequality, the interval $\left[\tau^{*}, \tau^{*}+\delta\right]$ splits into $I_{0}^{*}$ and $I_{1}^{*}$. The third inequality uses the relation $\mu\left(J_{i}\right)=\lambda_{i} \mu\left(I_{1}^{*}\right)$.

Let us now turn to the condition (3): for each $t$ with $\tau^{*} \leq t \leq \tau^{*}+\delta$, we have

$$
\begin{aligned}
& \left\|w^{*}(t)-v(t)\right\|^{w} \\
& \leq\left\|w^{*}\left(\tau^{*}\right)-v\left(\tau^{*}\right)\right\|^{w}+2 M(t-\tau) \\
& \leq \psi\left(\tau^{*}\right)+2 M \epsilon \\
& \leq \psi(t)+2 M \epsilon
\end{aligned}
$$

Step 6. Since $w^{*}$ belongs to $\mathfrak{M}$, $u^{*}$ is not a maximal element of $\mathfrak{M}$, which is a contradiction. Thus we conclude that $u^{*}$ is defined on the whole interval $[0, T]$.

Step 7. We have shown, in the above various steps, that, for each solution $v$ of differential inclusion (II) and each $\epsilon>0$, there exists a solution $u^{*}$ of the differential inclusion (I) such that $\left\|u^{*}(t)-v(t)\right\|^{w} \leq \psi(t)+2 M \epsilon$ for all $t \in[0, T]$ and hence

$$
\begin{aligned}
& \sup _{t \in[0, T]}\left\|u^{*}(t)-v(t)\right\|^{w} \\
& \leq \psi(t)+2 M \epsilon \\
& \leq \epsilon\left(2 M e^{T}+4\left(e^{T}-1\right)+2 M\right) .
\end{aligned}
$$

Since $\epsilon>0$ is arbitrary, this completes the proof of Theorem 1 .

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# STRONGLY GRADED RINGS WHICH ARE MAXIMAL ORDERS 

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Abstract.
Let $R=\oplus_{n \in \mathbb{Z}} R_{n}$ be a strongly graded ring of type $\mathbb{Z}$. In [6], it is shown that if $R_{0}$ is a maximal order, then so is $R$. We define a concept of $\mathbb{Z}$-invariant maximal order and show $R_{0}$ is a $\mathbb{Z}$-invariant maximal order if and only if $R$ is a maximal order. We provide examples of $R_{0}$ which are $\mathbb{Z}$-invariant maximal orders but not maximal orders.

1 Introduction Let $R=\oplus_{n \in \mathbb{Z}} R_{n}$ be a strongly graded ring of type $\mathbb{Z}$, where $\mathbb{Z}$ is the ring of integers. We always assume that $R_{0}$, the degree zero part, is a prime Goldie ring with its quotient ring $Q_{0}$ and $C_{0}=\left\{c \in R_{0} \mid c\right.$ is regular in $\left.R_{0}\right\}$, which is a regular Ore set of $R$ and the ring of fractions $Q^{g}$ of $R$ at $C_{0}$ has the following properties:
(i). $Q^{g}=\oplus_{n \in \mathbb{Z}} Q_{0} R_{n}\left(Q_{0} R_{n}=R_{n} Q_{0}\right)$.
(ii). $Q^{g}=Q_{0}\left[X, X^{-1}, \sigma\right]$ for some automorphism $\sigma$ of $R_{0}([6,1.3])$ and so
(iii). $Q^{g}$ is a left and right principal ideal ring.

We denote by the quotient ring of $R$ by $Q$. We define a concept of $\mathbb{Z}$-invariant maximal order in order to get the following three conditions are equivalent: (i) $R_{0}$ is a $\mathbb{Z}$-invariant maximal order (ii) $R$ is a maximal order (iii) $R$ is a graded maximal order. We give examples of $R_{0}$ which are $\mathbb{Z}$-invariant maximal orders but not maximal orders. We refer the readers to [7] or [8] and [9] for some elementary properties and some definitions of order theory and graded ring theory which are not mentioned in the paper.

2 The proof of Theorem Since $Q^{g}$ is the quotient ring of $R$ at $C_{0}$, the following lemma follows from the proof of [2, Theorem 1.31].

Lemma 1 Let $A$ be an ideal of $R$. Then $A Q^{g}=Q^{g} A$.

Lemma 2 Let $A_{0}$ be an ideal of $R_{0}$. Then the right ideal $A_{0} R$ is an ideal of $R$ if and only if $R_{n} A_{0}=A_{0} R_{n}$ for all $n \in \mathbb{Z}$. In this case, $A_{0} R$ is a graded ideal.

Proof. If $A_{0} R$ is an ideal of $R$, then $R_{n} A_{0} R_{-n} \subseteq A_{0}$, that is, $R_{n} A_{0} \subseteq A_{0} R_{n}$ for all $n \in \mathbb{Z}$ and so $R_{-n} A_{0} \subseteq A_{0} R_{-n}$ also follows. Hence $R_{n} A_{0}=A_{0} R_{n}$. Conversely if $R_{n} A_{0}=A_{0} R_{n}$ for all $n \in \mathbb{Z}$, then it is easy to see that $A_{0} R$ is an ideal of $R$.

[^2]
## Definition 1

(1). A left and right $R_{0}$-submodule $A_{0}$ of $Q_{0}$ is called $\mathbb{Z}$-invariant if $R_{n} A_{0}=A_{0} R_{n}$ for all $n \in \mathbb{Z}$.
(2). $R_{0}$ is called a $\mathbb{Z}$-invariant maximal order in $Q_{0}$ if $O_{l}\left(A_{0}\right)=R_{0}=O_{r}\left(A_{0}\right)$ for any nonzero $\mathbb{Z}$-invariant ideal $A_{0}$ of $R_{0}$.
(3). (10, p.205) $R$ is a graded maximal order in $Q^{g}$ if for each graded over-ring $S$ such that $R \subseteq S \subseteq Q^{g}$ and $a S b \subseteq R$ for some regular homogeneous elements $a, b \in Q^{g}$, it follows $R=S$.

## Lemma 3

(1). Let $A_{0}$ and $B_{0}$ be $\mathbb{Z}$-invariant left and right $R_{0}$-submodules in $Q_{0}$. Then $A_{0} B_{0}$ is $\mathbb{Z}$-invariant.
(2). Let $A_{0}$ be a $\mathbb{Z}$-invariant left $R_{0}$-ideal which is a right $R_{0}$-submodule in $Q_{0}$ and $B_{0}$ be a $\mathbb{Z}$-invariant right $R_{0}$-ideal which is a left $R_{0}$-submodule in $Q_{0}$. Then $C_{0}=\left\{r_{0} \in R_{0} \mid\right.$ $\left.A_{0} r_{0} \subseteq R_{0}\right\}$ and $D_{0}=\left\{r_{0} \in R_{0} \mid r_{0} B_{0} \subseteq R_{0}\right\}$ ) are both $\mathbb{Z}$-invariant.

Proof.
(1). It is clear.
(2). $R_{0} \supseteq A_{0} C_{0}$ implies $R_{0} \supseteq R_{n} A_{0} C_{0} R_{-n}=A_{0} R_{n} C_{0} R_{-n}$ for all $n \in \mathbb{Z}$ and so $R_{n} C_{0} R_{-n} \subseteq$ $C_{0}$ and also $R_{-n} C_{0} R_{n} \subseteq C_{0}$. Hence $C_{0} R_{n}=R_{n} C_{0}$ for all $n \in \mathbb{Z}$, that is, $C_{0}$ is $\mathbb{Z}^{-}$ invariant. Similarly $D_{0}$ is $\mathbb{Z}$-invariant.

Lemma 4 The following conditions are equivalent.
(1). $R_{0}$ is a $\mathbb{Z}$-invariant maximal order.
(2). $O_{l}\left(A_{0}\right)=R_{0}$ for each $\mathbb{Z}$-invariant left $R_{0}$-ideal $A_{0}$ which is a right $R_{0}$-submodule in $Q_{0}$, and $O_{r}\left(B_{0}\right)=R_{0}$ for each $\mathbb{Z}$-invariant right $R_{0}$-ideal $B_{0}$ which is a left $R_{0}$-submodule in $Q_{0}$.

Proof.
$(2) \Rightarrow(1)$ This is a special case.
$(1) \Rightarrow(2)$ Let $A_{0}$ be a $\mathbb{Z}$-invariant left $R_{0}$-ideal which is a right $R_{0}$-submodule in $Q_{0}$ and let $C_{0}=\left\{r_{0} \in R_{0} \mid A_{0} r_{0} \subseteq R_{0}\right\}$. Then $A_{0} C_{0}$ is a $\mathbb{Z}$-invariant ideal of $R_{0}$ by Lemma 3 . Thus $R_{0}=O_{l}\left(A_{0} C_{0}\right) \supseteq O_{l}\left(A_{0}\right) \supseteq R_{0}$ and so $O_{l}\left(A_{0}\right)=R_{0}$ follows. Similarly if $B_{0}$ is a $\mathbb{Z}$-invariant right $R_{0}$-ideal which is a left $R_{0}$-submodule in $Q_{0}$, then $O_{r}\left(B_{0}\right)=R_{0}$.

Theorem 1 Let $R=\oplus_{n \in \mathbb{Z}} R_{n}$ be a strongly graded ring of type $\mathbb{Z}$. Then the following conditions are equivalent:
(1). $R_{0}$ is a $\mathbb{Z}$-invariant maximal order in $Q_{0}$.
(2). $R$ is a maximal order in $Q$.
(3). $R$ is a graded maximal order in $Q^{g}$.

Proof.
(1) $\Rightarrow(2)$ Let $S$ be an over-ring of $R$ such that $a S b \subseteq R$ for some regular $a, b \in Q$. We may assume that $a, b \in R$. Put $T=R+R a S$, an over-ring of $R$ with $T b \subseteq R$. We claim $T=R$. Since $T b R$ is an ideal of $R$, it follows from Lemma 1 that $T b Q^{g}=T b R Q^{g}=$ $u Q^{g}=Q^{g} u$ for some regular element $u \in Q^{g}$ since $Q^{g}$ is a principal ideal ring. For any $t \in T$, tu $\in T b Q^{g}=Q^{g} u$ and so $t \in Q^{g}$. Thus $T \subseteq Q^{g}$ follows. For any $n \in \mathbb{Z}$, let $C_{n}(T)=\left\{a_{n} \in Q_{0} R_{n} \mid \exists t=a_{n}+a_{n_{1}}+\cdots+a_{n_{l}} \in T\right.$ such that $\left.n>n_{i}(1 \leq i \leq l)\right\} \cup\{0\}$, which is a left and right $R_{0}$-submodule of $Q^{g}$. It is easy to see that $C_{n}(T)=R_{n} C_{0}(T)=$ $C_{0}(T) R_{n}$. So, in particular, $C_{0}(T)$ is a $\mathbb{Z}$-invariant over-ring of $R_{0}$. To prove that $C_{0}(T)$ is a left $R_{0}$-ideal, write $b=b_{k}+$ (the lower degree parts). Since $T b \subseteq R$, it follows that $R_{0} \supseteq C_{-k}(T) b_{k}=C_{0}(T) R_{-k} b_{k}$ and so $R_{0} \supseteq C_{0}(T) R_{-k} b_{k} R_{0}$. Hence $C_{0}(T)$ is a left $R_{0}$-ideal since $R_{-k} b_{k} R_{0}$ is a non-zero ideal of $R_{0}$ and is a right $R_{0}$-submodule. Thus, by Lemma $4, O_{l}\left(C_{0}(T)\right)=R_{0}$ and $R_{0} \subseteq C_{0}(T) \subseteq O_{l}\left(C_{0}(T)\right)=R_{0}$ since $C_{0}(T)$ is an over-ring of $R_{0}$, which implies $R_{0}=C_{0}(T)$ and $R_{n}=R_{n} C_{0}(T)=C_{n}(T)$ for all $n \in \mathbb{Z}$. Hence $T=R$ follows. Since $a S \subseteq R a S \subseteq T=R$, the left version of the above proof shows that $S=R$. Hence $R$ is a maximal order in $Q$.
$(2) \Rightarrow(3)$ This is a special case.
(3) $\Rightarrow(1)$ Let $A_{0}$ be a $\mathbb{Z}$-invariant ideal of $R_{0}$. By Lemma 2, $A_{0} R$ is a graded ideal of $R$. Thus it follows from [6, Lemma 1.5] that $R O_{l}\left(A_{0}\right)=O_{l}\left(A_{0} R\right)=R=O_{r}\left(R A_{0}\right)=O_{r}\left(A_{0}\right) R$ and so $O_{l}\left(A_{0}\right)=R_{0}=O_{r}\left(A_{0}\right)$. Hence $R_{0}$ is a $\mathbb{Z}$-invariant maximal order in $Q_{0}$.

Finally, we give some examples of maximal orders $R$ such that $R_{0}$ are $\mathbb{Z}$-invariant maximal orders but not maximal orders.

Let $R_{0}$ be a hereditary Noetherian prime ring (an HNP ring for short) with its quotient ring $Q_{0}$ satisfying the following conditions:
(a). There is a cycle $M_{01}, \ldots, M_{0 n}(n \geq 2)$ so that $X=M_{01} \cap \cdots \cap M_{0 n}$ is a maximal invertible ideal of $R_{0}$.
(b). Any maximal ideal different from $M_{0 i}(1 \leq i \leq n)$ is invertible.

See [1] and [5] for examples of HNP rings satisfying (a) and (b) ( the simplest example is $\left(\begin{array}{ll}\mathbb{Z} & p \mathbb{Z} \\ \mathbb{Z} & \mathbb{Z}\end{array}\right)$, where $p$ is a prime number). Let

$$
R=\oplus_{n \in \mathbb{Z}} X^{n}\left(R_{n}=X^{n}\right)
$$

a strongly graded ring of type $\mathbb{Z}$, and $A_{0}$ be an eventually idempotent ideal of $R_{0}$. Then there are $M_{0 i_{1}}, \ldots, M_{0 i_{r}} i_{j} \in\{1, \ldots, n\}(r<n)$ which are the full set of maximal ideals containing $A_{0}$. Thus $A_{0}$ is not a $\mathbb{Z}$-invariant ideal by [4, Theorem 14].

Hence $R_{0}$ is a $\mathbb{Z}$-invariant maximal order since an ideal of $R_{0}$ is $\mathbb{Z}$-invariant if and only if it is invertible by [3, Theorem 2.9 and 4.2].

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# VERMA MODULES OVER $A \mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ GRADED SUPERALGEBRA AND INVARIANT DIFFERENTIAL EQUATIONS 

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#### Abstract

Lowest weight representations of the $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded superalgebra introduced by Rittenberg and Wyler are investigated. We give a explicit construction of Verma modules over the $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded superalgebra and show their reducibility by using singular vectors. The explicit formula of singular vectors are given and are used to derive partial differential equations invariant under the color supergroup generated by the $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded superalgebra.


1 Introduction. The present work aims to study representations of a $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded Lie superalgebra (also called color superalgebra) and its application to differential equations. We focus on a simple example of color superalgebras for the sake of simplicity and investigate Verma modules over it. The result is used to derive the partial differential equations which are invariant under the $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded Lie supergroup generated by the color superalgebra.

Color superalgebras are a generalization of Lie superalgebras introduced by Rittenberg and Wyler [1, 2] (see also [3, 4]). The idea of generalization is to extend the $\mathbb{Z}_{2}$ graded structure of the underlying vector space of Lie superalgebra to more general abelian groups. The group $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$, discussed here, is the simplest non-trivial example of the generalization. During the last four decades, structure theory, classification of possible color superalgebras of given dimension etc. have been studied by many authors. See, for example, $[5,6,7,8,9$, $10,11,12,13$ ] and references therein. However, applications of color superalgebras to mathematical and physical problems are very limited [ $14,15,16,17,18,19,20,21,22,23]$. In $[22,23]$ it is shown that symmetries of a first order linear partial differential equation, called the Lévy-Leblond equation [24], are generated by a $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded superalgebra. The authors seek symmetry operators of the Lévy-Leblond equation systematically and found that they do not close as a Lie algebra or superalgebra but do in a $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded superalgebra. Here, we reverse the argument. We start with a simple example of $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded superalgebra and then study its lowest weight representations and their irreducibility. Reducibility of the representations is detected by the existence of singular vectors. Explicit formulae of the singular vectors allows us to write down invariant partial differential equations. This is a generalization of the method developed for semi-simple Lie groups [25] to the $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ setting. Thus we shall see the method is valid beyond the semi-simple Lie groups.

The plan of this paper is as follows. In the next section we give a definition of $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded superalgebra and present the one, denoted by $\mathfrak{g}$, investigated in this work. Lowest weight Verma modules over $\mathfrak{g}$ are studied in $\S 3$. After discussing the failure of the naive approach, we construct Verma modules over $\mathfrak{g}$ by using a trick. A list of all singular vectors in the Verma modules is presented with an explicit formulae of the singular vectors. The formulae of the singular vectors are used to find differential equations whose symmetries are generated by $\mathfrak{g}$ in $\S 4$. Summary and some remarks are given in $\S 5$.

[^3]2 Definition of $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded superalgebra. Let $\mathfrak{g}$ be a vector space over $\mathbb{C}$ (or $\mathbb{R}$ ) which is a direct sum of four subspaces labelled by an element of the group $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ :

$$
\begin{equation*}
\mathfrak{g}=\mathfrak{g}_{(0,0)} \oplus \mathfrak{g}_{(0,1)} \oplus \mathfrak{g}_{(1,0)} \oplus \mathfrak{g}_{(1,1)} \tag{2.1}
\end{equation*}
$$

An inner product of two $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ vectors $\boldsymbol{\alpha}=\left(\alpha_{1}, \alpha_{2}\right), \boldsymbol{\beta}=\left(\beta_{1}, \beta_{2}\right)$ is defined as usual:

$$
\begin{equation*}
\boldsymbol{\alpha} \cdot \boldsymbol{\beta}=\alpha_{1} \beta_{1}+\alpha_{2} \beta_{2} \tag{2.2}
\end{equation*}
$$

Now we give a definition of $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded superalgebra according to [1, 2].
Definition 2.1. If $\mathfrak{g}$ admits a bilinear form $\llbracket, \rrbracket: \mathfrak{g} \times \mathfrak{g} \rightarrow \mathfrak{g}$ satisfying the following three relations, then $\mathfrak{g}$ is called a $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded color superalgebra:

1. $\llbracket \mathfrak{g}_{\boldsymbol{\alpha}}, \mathfrak{g}_{\boldsymbol{\beta}} \rrbracket \subseteq \mathfrak{g}_{\boldsymbol{\alpha}+\boldsymbol{\beta}}$,
2. $\llbracket X_{\boldsymbol{\alpha}}, X_{\boldsymbol{\beta}} \rrbracket=-(-1)^{\boldsymbol{\alpha} \cdot \boldsymbol{\beta}} \llbracket X_{\boldsymbol{\beta}}, X_{\boldsymbol{\alpha}} \rrbracket$,
3. $\llbracket X_{\boldsymbol{\alpha}}, \llbracket X_{\boldsymbol{\beta}}, X_{\boldsymbol{\gamma}} \rrbracket \rrbracket(-1)^{\boldsymbol{\alpha} \cdot \gamma}+$ cyclic perm. $=0$,
where $X_{\boldsymbol{\alpha}} \in \mathfrak{g}_{\boldsymbol{\alpha}}$ and the third relation is called the graded Jacobi identity.
When the inner product $\boldsymbol{\alpha} \cdot \boldsymbol{\beta}$ is an even integer the graded Lie bracket $\llbracket, \rrbracket$ is understood as a commutator, while it is an anticommutator if $\boldsymbol{\alpha} \cdot \boldsymbol{\beta}$ is an odd integer. If $N=1$, then Definition 2.1 is identical to the definition of Lie superalgebras so that the color superalgebra is a natural generalization of Lie superalgebra.

The $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded superalgebra investigated in this work is eight dimensional and its basis and grading are given as follows:

$$
\begin{array}{lll}
(0,0) & : & A_{-}, A_{+}, N \\
(1,0) & : & b_{-}, b_{+} \\
(0,1) & : & a_{-}, a_{+}  \tag{2.3}\\
(1,1) & : & F
\end{array}
$$

The basis satisfy the relations (only non-vanishing ones are presented):

$$
\begin{align*}
& {\left[A_{-}, A_{+}\right]=4 N, \quad\left[A_{-}, N\right]=2 A_{-}, \quad\left[A_{+}, N\right]=-2 A_{+},} \\
& {\left[A_{-}, b_{+}\right]=2 b_{-}, \quad\left[A_{+}, b_{-}\right]=-2 b_{+}, \quad\left[N, b_{-}\right]=-b_{-}, \quad\left[N, b_{+}\right]=b_{+},} \\
& {\left[A_{-}, a_{+}\right]=2 a_{-}, \quad\left[A_{+}, a_{-}\right]=-2 a_{+}, \quad\left[N, a_{-}\right]=-a_{-}, \quad\left[N, a_{+}\right]=a_{+},} \\
& \left\{b_{-}, b_{-}\right\}=2 A_{-}, \quad\left\{b_{-}, b_{+}\right\}=2 N, \quad\left\{b_{+}, b_{+}\right\}=2 A_{+}, \\
& {\left[b_{-}, a_{+}\right]=F, \quad\left[b_{+}, a_{-}\right]=-F, \quad\left\{b_{-}, F\right\}=2 a_{-}, \quad\left\{b_{+}, F\right\}=2 a_{+},} \\
& \left\{a_{-}, a_{-}\right\}=2 A_{-}, \quad\left\{a_{-}, a_{+}\right\}=2 N, \quad\left\{a_{+}, a_{+}\right\}=2 A_{+}, \\
& \left\{a_{-}, F\right\}=2 b_{-}, \quad\left\{a_{+}, F\right\}=2 b_{+} . \tag{2.4}
\end{align*}
$$

One may see from the defining relations that the color superalgebra has two osp $(1,2)$ subalgebras, $\left\langle A_{ \pm}, N, a_{ \pm}\right\rangle$and $\left\langle A_{ \pm}, N, b_{ \pm}\right\rangle$. This color superalgebra was given in [1] as one of the non-trivial examples of $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded superalgebras. It is also discussed in [13] as an $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ extension of the superalgebra generated by boson and fermion operators. We denote the color superalgebra (2.3) by $\mathfrak{g}$.

Lemma 2.2. The superalgebra $\mathfrak{g}$ admits an algebra anti-involution $\omega: \mathfrak{g} \rightarrow \mathfrak{g}$ defined by

$$
\begin{equation*}
\omega\left(X_{ \pm}\right)=X_{\mp}, \quad \omega(Y)=Y, \quad X_{ \pm}=A_{ \pm}, a_{ \pm}, b_{ \pm}, \quad Y=N, F \tag{2.5}
\end{equation*}
$$

Proof. It is checked by straightforward computation.
Next we give a definition of the $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded Grassmann numbers which was also introduced in [2].

Definition 2.3. Let $\boldsymbol{\alpha} \in \mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ and $\zeta_{\boldsymbol{\alpha}, i}$ be a basis of a vector space over the same field as $\mathfrak{g}$. If the basis satisfies the relations

$$
\begin{equation*}
\llbracket \zeta_{\boldsymbol{\alpha}, i}, \zeta_{\boldsymbol{\beta}, j} \rrbracket=0 \tag{2.6}
\end{equation*}
$$

then we call $\zeta_{\alpha, i}$ the $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded Grassmann numbers.
An extension of an integral and a derivative of the ordinary Grassmann numbers to $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ setting is discussed in [11]. We use the $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded Grassmann numbers in the subsequent sections. The basis $\zeta_{\boldsymbol{\alpha}, i}$ may be realized in terms of the ordinary Grassmann numbers and the Clifford algebra [13]. Let us recall that the Clifford algebra $C l(p, q)$ is a unital algebra generated by $\gamma_{i}(i=1,2, \ldots, N=p+q)$ subject to the relations:

$$
\begin{equation*}
\left\{\gamma_{i}, \gamma_{j}\right\}=2 \eta_{i j}, \quad \eta=\operatorname{diag}(\underbrace{+1, \ldots,+1}_{p}, \underbrace{-1, \ldots,-1}_{q}) \tag{2.7}
\end{equation*}
$$

The $C l(p, q)$ is a $2^{N}$ dimensional algebra whose elements are given by a product of the generators:

$$
1, \gamma_{i}, \gamma_{i} \gamma_{j}, \gamma_{i} \gamma_{j} \gamma_{k}, \ldots, \gamma_{1} \gamma_{2} \cdots \gamma_{N}
$$

Lemma 2.4. The $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded Grassmann numbers $\zeta_{\boldsymbol{\alpha}, i}$ are realized in terms of the ordinary Grassmann number $\xi_{\mu}$ and the Clifford algebra $C l(p, q), p+q=2$ :

$$
\begin{aligned}
& \zeta_{(0,0), m}=1 \otimes x_{m}, \quad \zeta_{(1,0), \mu}=\gamma_{1} \otimes \xi_{\mu} \\
& \zeta_{(0,1), \mu}=\gamma_{2} \otimes \xi_{\mu}, \quad \zeta_{(1,1), m}=\gamma_{1} \gamma_{2} \otimes x_{m}
\end{aligned}
$$

where $x_{m} \in \mathbb{R}$.
3 Verma modules and their reducibility. We want to investigate the lowest weight representations of $\mathfrak{g}$ by employing the standard procedure of Lie theory [26]. However, it turns out that the naive approach is not successful in defining Verma modules over $\mathfrak{g}$. We begin with the failure of the naive construction of Verma modules and then consider a modified approach.
3.1 Failure of the naive construction. The algebra $\mathfrak{g}$ has a natural triangular decomposition:

$$
\begin{align*}
\mathfrak{g}_{+} & =\left\langle A_{+}, a_{+}, b_{+}\right\rangle \\
\mathfrak{g}_{0} & =\langle N, F\rangle \\
\mathfrak{g}_{-} & =\left\langle A_{-}, a_{-}, b_{-}\right\rangle \tag{3.1}
\end{align*}
$$

This is based on the eigenvalues of $\operatorname{ad} N$ which are given as follows:

$$
\begin{array}{ccccc}
+2 & +1 & 0 & -1 & -2 \\
\hline A_{+} & a_{+}, b_{+} & N, F & a_{-}, b_{-} & A_{-}
\end{array}
$$

The anti-involution $\omega$ (2.5) acts on the subspaces as

$$
\begin{equation*}
\omega\left(\mathfrak{g}_{ \pm}\right)=\mathfrak{g}_{\mp}, \quad \omega\left(\mathfrak{g}_{0}\right)=\mathfrak{g}_{0} \tag{3.2}
\end{equation*}
$$

and it is easy to see that $\llbracket \mathfrak{g}_{0}, \mathfrak{g}_{ \pm} \rrbracket \subseteq \mathfrak{g}_{ \pm}$and $\llbracket \mathfrak{g}_{ \pm}, \mathfrak{g}_{ \pm} \rrbracket \subseteq \mathfrak{g}_{ \pm}$. Therefore, (3.1) is a natural triangular decomposition of $\mathfrak{g}$.

The decomposition (3.1) leads us to define Verma modules in a standard way. We first define the lowest weight state $|h, \varphi\rangle$ by

$$
\begin{equation*}
a_{-}|h, \varphi\rangle=0, \quad N|h, \varphi\rangle=h|h, \varphi\rangle, \quad F|h, \varphi\rangle=\varphi|h, \varphi\rangle \tag{3.3}
\end{equation*}
$$

It follows that

$$
\begin{equation*}
A_{-}|h, \varphi\rangle=b_{-}|h, \varphi\rangle=0 \tag{3.4}
\end{equation*}
$$

Then the Verma modules over $\mathfrak{g}$ are defined by $M(h, \varphi)=\mathcal{U}\left(\mathfrak{g}_{+}\right) \otimes|h, \varphi\rangle$ where $\mathcal{U}\left(\mathfrak{g}_{+}\right)$is the universal enveloping algebra of $\mathfrak{g}_{+}$. The basis of $M(h, \varphi)$ is obviously given by

$$
\begin{equation*}
|k, \ell\rangle=\left(a_{+}\right)^{k}\left(b_{+}\right)^{\ell}|h, \varphi\rangle, \quad k, \ell \in \mathbb{Z}_{\geq 0} \tag{3.5}
\end{equation*}
$$

It is immediate to see that the relations $\left\{a_{+}, a_{+}\right\}=\left\{b_{+}, b_{+}\right\}=2 A_{+}$are not realized on this basis. This shows the failure of the naive construction of Verma modules.
3.2 Verma modules and singular vectors. To overcome the difficulty in $\S 3.1$ we use a $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded Grassmann number $\zeta$ of degree $(1,1)$. Suppose further that $\zeta^{2}=1$. Legitimacy of the assumption is ensured by Lemma 2.4 since $\zeta$ may be realized by using $C l(1,1)$. Now we define a new basis of $\mathfrak{g}$ :

$$
\begin{equation*}
c_{ \pm}=\frac{1}{\sqrt{2}}\left(a_{-} \pm b_{-} \zeta\right), \quad d_{ \pm}=\frac{1}{\sqrt{2}}\left(a_{+} \pm \zeta b_{+}\right), \quad \tilde{F}=\frac{1}{2} \zeta F . \tag{3.6}
\end{equation*}
$$

The non-vanishing defining relations for the new basis are written as

$$
\begin{align*}
\left\{c_{+}, c_{-}\right\} & =2 A_{-}, & \left\{d_{+}, d_{-}\right\} & =2 A_{+},
\end{align*} \quad\left\{c_{ \pm}, d_{ \pm}\right\}=2(N \mp \tilde{F}),
$$

We remark that $c_{ \pm}$and $d_{ \pm}$are nilpotent and that $c_{ \pm}$anticommutes with $d_{\mp}$. The antiinvolution $\omega(2.5)$ is extended to the new basis (3.6) by setting $\omega(\zeta)=1$ :

$$
\begin{equation*}
\omega\left(c_{ \pm}\right)=d_{ \pm}, \quad \omega(\tilde{F})=\tilde{F} \tag{3.8}
\end{equation*}
$$

In fact, this change of basis converts the $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded superalgebra to an ordinary superalgebra of $\mathbb{Z}_{2}$ grading. Due to the degree of $\zeta$, the degree of $c_{ \pm}$and $d_{ \pm}$are all $(0,1)$, while the degree of $\tilde{F}$ is $(0,0)$. Thus one may consider representations of the $\mathbb{Z}_{2}$ graded algebra, then convert it to the ones for $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ grading.

The present choice of the basis of $\mathfrak{g}$ diagonalize $\operatorname{ad} N$ and $\operatorname{ad} \tilde{F}$. Their eigenvalues are summarized as

$$
\begin{array}{llll}
A_{+}(+2,0), & A_{-}(-2,0), & c_{+}(-1,-1), & c_{-}(-1,+1) \\
N(0,0), & \tilde{F}(0,0), & d_{+}(+1,+1), & d_{-}(+1,-1) \tag{3.9}
\end{array}
$$

We introduce the triangular decomposition of $\mathfrak{g}$ according to the eigenvalue of $\operatorname{ad} N$ :

$$
\begin{equation*}
\mathfrak{g}_{+}=\left\langle A_{+}, d_{ \pm}\right\rangle, \quad \mathfrak{g}_{0}=\langle N, \tilde{F}\rangle, \quad \mathfrak{g}_{-}=\left\langle A_{-}, c_{ \pm}\right\rangle \tag{3.10}
\end{equation*}
$$

Define the lowest weight vector $|h, f\rangle$ by

$$
\begin{equation*}
c_{ \pm}|h, f\rangle=0, \quad N|h, f\rangle=h|h, f\rangle, \quad \tilde{F}|h, f\rangle=f|h, f\rangle \tag{3.11}
\end{equation*}
$$

It then follows that $A_{-}|h, f\rangle=0$. We define the Verma modules over $\mathfrak{g}$ by a space induced from $|h, f\rangle$ by $M(h, f)=\mathcal{U}\left(\mathfrak{g}_{+}\right) \otimes|h, f\rangle$. The natural basis of $M(h, f)$ is given by

$$
\begin{equation*}
|k, \mu, \nu\rangle=\left(A_{+}\right)^{k}\left(d_{+}\right)^{\mu}\left(d_{-}\right)^{\nu}|h, f\rangle, \quad k \in \mathbb{Z}_{\geq 0}, \mu, \nu \in\{0,1\} \tag{3.12}
\end{equation*}
$$

It is then straightforward to compute the action of $\mathfrak{g}$ on $M(h, f)$. The action of $\mathfrak{g}_{0}$ yields

$$
\begin{align*}
N|k, \mu, \nu\rangle & =(h+2 k+\mu+\nu)|k, \mu, \nu\rangle \\
\tilde{F}|k, \mu, \nu\rangle & =(f+\mu-\nu)|k, \mu, \nu\rangle \tag{3.13}
\end{align*}
$$

The action of $\mathfrak{g}_{+}$is given by

$$
\begin{align*}
A_{+}|k, \mu, \nu\rangle & =|k+1, \mu, \nu\rangle \\
d_{+}|k, \mu, \nu\rangle & =\delta_{\mu, 0}|k, \mu+1, \nu\rangle \\
d_{-}|k, \mu, \nu\rangle & =(-1)^{\mu} \delta_{\nu, 0}|k, \mu, \nu+1\rangle+\delta_{\mu, 1} 2|k+1, \mu-1, \nu\rangle . \tag{3.14}
\end{align*}
$$

The action of $\mathfrak{g}_{-}$is summarized as

$$
\begin{align*}
A_{-}|k, \mu, \nu\rangle & =4 k(h+k+\mu+\nu-1)|k-1, \mu, \nu\rangle+\delta_{\mu, 1} \delta_{\nu, 1} 4(h+f)|k, \mu-1, \nu-1\rangle \\
c+|k, \mu, \nu\rangle & =\delta_{\mu, 1} 2(h+2 k+2 \nu-f)|k, \mu-1, \nu\rangle+(-1)^{\mu} \delta_{\nu, 0} 2 k|k-1, \mu, \nu+1\rangle \\
c_{-}|k, \mu, \nu\rangle & =(-1)^{\mu} \delta_{\nu, 1} 2(h+f)|k, \mu, \nu-1\rangle+\delta_{\mu, 0} 2 k|k-1, \mu+1, \nu\rangle \tag{3.15}
\end{align*}
$$

We have successfully obtained the Verma modules over $\mathfrak{g}$.
Let us now discuss reducibility of the Verma modules. This may be done by singular vectors [26]. The existence of the singular vector, by definition, means that $M(h, f)$ is reducible. The Verma module $M(h, f)$ has a natural grading as a vector space:

$$
\begin{equation*}
\left.M(h, f)=\bigoplus_{m \in \mathbb{N}} M_{m}(h, f), \quad M_{m}(h, f)=\{|v\rangle \in M(h, f)|N| v\rangle=(h+m)|v\rangle\right\} \tag{3.16}
\end{equation*}
$$

where $\mathbb{N}$ is the set of non-negative integers. In the basis (3.12), the integer $m$ is given by $m=2 k+\mu+\nu$ and we call $m$ the level. Any singular vector is an eigenvector of $N$ so that it must be an element of a subspace $M_{m}(h, f)$. We give a complete list of singular vectors in $M(h, f)$.

Theorem 3.1. $M(h, f)$ has precisely one singular vector if $h, f$ satisfy one of the following conditions:
(1) $h=f$ : The singular vector is $|0,1,0\rangle$.
(2) $h=-f$ : The singular vector is $|0,0,1\rangle$.
(3) $h=-n$ and $f \neq n$ for a positive integer $n$ : The singular vector exists at level $2 n$ subspace and given by

$$
\begin{equation*}
|n, 0,0\rangle+\frac{n}{f-n}|n-1,1,1\rangle \tag{3.17}
\end{equation*}
$$

Proof. First of all, we note that $\operatorname{dim} M_{m}(h, f)=2$ for any level $m$. If the level $m=2 n+1$ is odd, then $M_{m}(h, f)$ is spanned by $|n, 0,1\rangle$ and $|n, 1,0\rangle$, while if $m=2 n$, then $M_{m}(h, f)$ is spanned by $|n, 0,0\rangle$ and $|n-1,1,1\rangle$. We study the odd and even level separately.
(i) $m=2 n+1$. Since $|n, 0,1\rangle$ and $|n, 1,0\rangle$ have distinct eigenvalues of $\tilde{F}$, singular vector is not a linear combination of the two vectors. Requirement that $|n, 0,1\rangle$ is annihilated by the action of $\mathfrak{g}_{-}$provides the relations

$$
\begin{equation*}
n=0, \quad h+f=0 \tag{3.18}
\end{equation*}
$$

The same argument for $|n, 1,0\rangle$ provides the relations

$$
\begin{equation*}
n=0, \quad h-f=0 \tag{3.19}
\end{equation*}
$$

This proves (1) and (2) of the theorem.
(ii) $m=2 n$. A singular vector, if any, may be written as

$$
\begin{equation*}
|u\rangle=|n, 0,0\rangle+\alpha|n-1,1,1\rangle \tag{3.20}
\end{equation*}
$$

with a constant $\alpha .|u\rangle$ is an eigenvector of $\tilde{F}$ with the eigenvalue $f$. The requirement that the action of $\mathfrak{g}_{-}$annihilate $|u\rangle$ gives the relations

$$
\begin{align*}
& n+\alpha(h+2 n-f)=0 \\
& n-\alpha(h+f)=0 \\
& n(h+n-1)+\alpha(h+f)=0 \\
& \alpha(n-1)(h+n)=0 \tag{3.21}
\end{align*}
$$

The last equation implies that there exists three possibilities:
(a) $\alpha=0$,
(b) $h=-n$,
(c) $n=1$.

The case (a) means that $n=0$ so that $|u\rangle=|0,0,0\rangle$. Thus this case is trivial. The case (b), the first three equations in (3.21) are reduced to single relation:

$$
\begin{equation*}
n+\alpha(n-f)=0 \tag{3.22}
\end{equation*}
$$

If $f=n$, then the relations means that $n=0$. Thus we need $\alpha=0$ and we have a only trivial solution. If $f \neq n$, then (3.22) is solved to $\alpha$ and give the unique singular vector (3.17). The case (c) is reduced to a subclass of the case (b) since we find the relation $h=-1=-n$. We thus completed the proof.

Corollary 3.2. The Verma module $M(h, f)$ is irreducible if following conditions are true:
(1) $h \neq \pm f$
(2) $h \neq-n$ or $f=n$

4 Invariant partial differential equations. In this section, we derive partial differential equations which are invariant under the color supergroup generated by $\mathfrak{g}$. This is done by employing the method of [25]. We give a brief outline of the method and refer to [25] for detail. The basic idea of the method is to realize the Verma modules, discussed in $\S 3.2$, by a space of differentiable functions defined on the color supergroup.

Let $\mathcal{G}$ be the color supergroup generated by $\mathfrak{g}$ which means that an element $g \in \mathcal{G}$ is written as $g=\exp (X), X \in \mathfrak{g}$. We consider a space of functions $C_{\Lambda}$ on $\mathcal{G}$ having a special property. That is, $F \in C_{\Lambda}$ satisfies the relation:

$$
\begin{equation*}
F\left(g g_{0} g_{-}\right)=\exp (\Lambda(H)) F(g), \quad g \in \mathcal{G}, g_{0}=\exp (H), H \in \mathfrak{g}_{0}, g_{-} \in \exp \left(\mathfrak{g}_{-}\right) \tag{4.1}
\end{equation*}
$$

where $\Lambda(H)$ is an eigenvalue of $H$. This property means that $F$ is, in fact, a function on $\mathcal{G}_{+}=\exp \left(\mathfrak{g}_{+}\right)$. We then define the left $\left(\pi_{L}\right)$ and right $\left(\pi_{R}\right)$ actions of $\mathfrak{g}$ on $C_{\Lambda}$ according to the standard way of Lie theory:

$$
\begin{align*}
\pi_{L}(X) F(g) & =\left.\frac{d}{d \tau} F\left(e^{-\tau X} g\right)\right|_{\tau=0}  \tag{4.2}\\
\pi_{R}(X) F(g) & =\left.\frac{d}{d \tau} F\left(g e^{\tau X}\right)\right|_{\tau=0} \tag{4.3}
\end{align*}
$$

where $X \in \mathfrak{g}$. Note that the parameter $\tau$ is a $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded Grassmann number of degree same as $X$. A definition of derivative with respect to the $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded Grassmann numbers is found in [11]. It is then easy to see that, due to the property (4.1), the function $F \in C_{\Lambda}$ plays the role of a lowest weight vector $|h, f\rangle$ by the right action. Namely, $\pi_{R}(X) F=0$ if $X \in \mathfrak{g}_{-}$and $\pi_{R}(X) F=\Lambda(x) F$ if $X \in \mathfrak{g}_{0}$. While $\pi_{R}(X)$ with $X \in \mathfrak{g}_{+}$becomes a differential operator so that a Verma module $M(h, f)$ is realized by the action of $\pi_{R}(X)$ with $X \in \mathfrak{g}_{+}$ on the function $F(g)$.

To be more explicit, we parametrize $g_{+} \in \mathcal{G}_{+}$as

$$
\begin{equation*}
g_{+}=\exp \left(x A_{+}\right) \exp \left(\psi_{+} d_{+}\right) \exp \left(\psi_{-} d_{-}\right) \tag{4.4}
\end{equation*}
$$

Then we obtain from (4.3)

$$
\begin{align*}
\pi_{R}\left(A_{+}\right) & =\frac{\partial}{\partial x} \\
\pi_{R}\left(d_{+}\right) & =\frac{\partial}{\partial \psi_{+}}+2 \psi_{-} \frac{\partial}{\partial x} \\
\pi_{R}\left(d_{-}\right) & =\frac{\partial}{\partial \psi_{-}} \tag{4.5}
\end{align*}
$$

and by setting $\Lambda(N)=h, \Lambda(\tilde{F})=f$ we have

$$
\begin{equation*}
\pi_{R}(N)=h, \quad \pi_{R}(\tilde{F})=f \tag{4.6}
\end{equation*}
$$

As shown in $\S 3.2$ the singular vectors in $M(h, f)$ is written as $\mathcal{P}|h, f\rangle$ with $\mathcal{P} \in U\left(\mathfrak{g}_{+}\right)$. It is shown in [25] that differential equations invariant under the color supergroup $\mathcal{G}$ are given by $\pi_{R}(\mathcal{P}) \varphi=0$. The symmetry transformations are generated by the left action (4.2). This is due to the fact that a singular vector is an intertwining operator of two representations $M(h, f)$ and $M\left(h^{\prime}, f^{\prime}\right)$. According to Theorem 3.1 we found a hierarchy of invariant differential equations.

Proposition 4.1. The following equations are invariant under the color supergroup $\mathcal{G}$ :

$$
\begin{align*}
& \frac{\partial}{\partial \psi_{-}} \varphi\left(x, \psi_{ \pm}\right)=0 \\
& \left(\frac{\partial}{\partial \psi_{+}}+2 \psi_{-} \frac{\partial}{\partial x}\right) \varphi\left(x, \psi_{ \pm}\right)=0 \\
& {\left[\frac{\partial}{\partial x}+\frac{n}{f-n}\left(\frac{\partial}{\partial \psi_{+}}+2 \psi_{-} \frac{\partial}{\partial x}\right) \frac{\partial}{\partial \psi_{-}}\right]\left(\frac{\partial}{\partial x}\right)^{n-1} \varphi\left(x, \psi_{ \pm}\right)=0} \tag{4.7}
\end{align*}
$$

Up to here our analysis of the representations of $\mathfrak{g}$ is done by converting the basis of $\mathfrak{g}$ into a $\mathbb{Z}_{2}$ grading ones. Thus the invariant equations are written in terms of the variables
of degree $(0,0)$ and $(0,1)$. We now rewrite the invariant differential equations in a form also containing the variables of degree $(1,0)$ and $(1,1)$. Let $\psi$ and $\theta$ be $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded Grassmann number of degree $(0,1)$ and $(1,0)$, respectively. Taking the basis change (3.6) into account we set

$$
\begin{equation*}
\psi_{ \pm}=\frac{1}{\sqrt{2}}(\psi \pm \zeta \theta) \tag{4.8}
\end{equation*}
$$

and replace $\psi_{ \pm}$with $\psi$ and $\theta$. Then the independent variables of the invariant differential equations become $x, \psi$ and $\theta$ and the equations yield as follows:

$$
\begin{aligned}
& \left(\frac{\partial}{\partial \psi}+\zeta \frac{\partial}{\partial \theta}\right) \varphi=0 \\
& {\left[\frac{\partial}{\partial \psi}-\zeta \frac{\partial}{\partial \theta}+2(\psi-\zeta \theta) \frac{\partial}{\partial x}\right] \varphi=0} \\
& {\left[\frac{\partial}{\partial x}+\frac{n}{f-n}\left(\psi \frac{\partial}{\partial \psi}+\theta \frac{\partial}{\partial \theta}-\zeta\left(\psi \frac{\partial}{\partial \theta}+\theta \frac{\partial}{\partial \psi}\right)\right) \frac{\partial}{\partial x}-\frac{n \zeta}{f-n} \frac{\partial^{2}}{\partial \psi \partial \theta}\right]\left(\frac{\partial}{\partial x}\right)^{n-1} \varphi=0}
\end{aligned}
$$

5 Concluding remarks. We studied lowest weight representations of the $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded superalgebra given by Rittenberg and Wyler. Reducibility of the Verma modules over the $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded superalgebra is shown by explicit construction of singular vectors. As an application of the present scheme invariant partial differential equations defined on the space of functions whose variables are $\mathbb{Z}_{2} \otimes \mathbb{Z}_{2}$ graded Grassmann numbers are obtained. The present scheme may apply to other color superalgebras and one may find many related differential equations. Since differential equations are very basic objects in theoretical physics and many areas of mathematics, the present line of investigation will reveal a connection of color superalgebras with different areas of science.

We close this paper with two possible lines of further investigations on applications of color superalgebras. The first one is mathematical application. It is well known that representations of Lie algebras and Lie groups are closely related to orthogonal polynomials. We anticipate that representations of color algebras or groups also have such connection. The second one is physical application. One may use a nonlinear realization method to write down Lagrangians which are invariant under a color supergroup, since the method is originally defined for Lie groups then extended to supergroups. Color supergroups are natural generalization of supergroups, thus nonlinear realization methods will be generalized to any color supergroups. If this is done, one may find a dynamical system whose symmetry is governed by a color supergroup. This would open the way to physical applications of color supergroups.
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# CAUCHY'S THEOREM FOR B-ALGEBRAS 

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#### Abstract

In this paper, we establish the Cauchy's Theorem for B-algebras. We also present some implications of Lagrange's Theorem and Cauchy's Theorem for Balgebras. In particular, the concept of $\mathrm{B}_{p}$-algebras is introduced.


1 Introduction In [9], the notion of B-algebras was introduced by J. Neggers and H.S. Kim. A B-algebra is an algebra $(X ; *, 0)$ of type $(2,0)$ (that is, a nonempty set $X$ with a binary operation $*$ and a constant 0 ) satisfying the following axioms for all $x, y, z \in X:$ (I) $x * x=0$, (II) $x * 0=x$, (III) $(x * y) * z=x *(z *(0 * y))$. A B-algebra $(X ; *, 0)$ is commutative [9] if $x *(0 * y)=y *(0 * x)$ for all $x, y \in X$. In [10], J. Neggers and H.S. Kim introduced the notions of a subalgebra and normality of B-algebras and some of their properties are established. A nonempty subset $N$ of $X$ is called a subalgebra of $X$ if $x * y \in N$ for any $x, y \in N$. It is called normal in $X$ if for any $x * y, a * b \in N$ implies $(x * a) *(y * b) \in N$. A normal subset of $X$ is a subalgebra of $X$. There are several properties of B-algebras as established by some authors [1-12]. The following properties are used in this paper, for any $x, y, z \in X$, we have (P1) $0 *(0 * x)=x[9]$, (P2) $x * y=0 *(y * x)$ [11], (P3) $x *(y * z)=(x *(0 * z)) * y[9],(\mathrm{P} 4) x * y=x * z$ implies $y=z[3],(\mathrm{P} 5)(0 * x) *(y * x)=0 * y$ [9]. In [2], J.S. Bantug and J.C. Endam established the Lagrange's Theorem for B-algebras. In this paper, we provide some partial results on the converse of this theorem. In particular, we establish the Cauchy's Theorem for B-algebras. As a consequence, we also introduce the concept of $\mathrm{B}_{p}$-algebras. Throughout this paper, $X$ means a B-algebra $(X ; *, 0)$.

2 Preliminaries This section presents some concepts and results needed in this paper. We start with some examples of B-algebras.

Example 2.1. [9] Let $X=\{0,1,2\}$ be a set with the following table of operation:

| $*$ | 0 | 1 | 2 |
| :--- | :--- | :--- | :--- |
| 0 | 0 | 2 | 1 |
| 1 | 1 | 0 | 2 |
| 2 | 2 | 1 | 0 |

Example 2.2. [9] Let $X=\{0,1,2,3,4,5\}$ be a set with the following table of operation:

| $*$ | 0 | 1 | 2 | 3 | 4 | 5 |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| 0 | 0 | 2 | 1 | 3 | 4 | 5 |
| 1 | 1 | 0 | 2 | 4 | 5 | 3 |
| 2 | 2 | 1 | 0 | 5 | 3 | 4 |
| 3 | 3 | 4 | 5 | 0 | 2 | 1 |
| 4 | 4 | 5 | 3 | 1 | 0 | 2 |
| 5 | 5 | 3 | 4 | 2 | 1 | 0 |

In [7], if $S$ is a subset of $X$, then $\langle S\rangle_{B}$ is the intersection of all subalgebra $H$ of $X$ such that $S \subseteq H$, and the subalgebra $\langle S\rangle_{B}$ of $X$ is called the subalgebra generated by $S$. If $X=\langle S\rangle_{B}$, then $S$ is called a set of generators for $X$. Moreover, $\langle S\rangle_{B}$ is the smallest subalgebra of $X$ containing $S$. If either $S=\varnothing$ or $S=\{0\}$, then $\langle S\rangle_{B}=\{0\}$. If $S$ is a subalgebra of $X$, then $\langle S\rangle_{B}=S$. In particular, $\langle X\rangle_{B}=X$.

Let $x \in X$. In [9], J. Neggers and H.S. Kim defined $x^{n}=x^{n-1} *(0 * x)$ for $n \geq 1$ and $x^{0}=0$. Then $x^{m} * x^{n}=x^{m-n}$ if $m \geq n$ and $x^{m} * x^{n}=0 * x^{n-m}$ otherwise. In [7], for each $x \in X$, N.C. Gonzaga and J.P. Vilela defined $-x=0 * x$ and $x^{-n}=(-x)^{n}$ for each $n \geq 1$. In [5], J.C. Endam and R.C. Teves defined $x^{m}=0 * x^{-m}$ for $m \leq-1$. If $m \geq 1$, then $x^{m}=0 *\left(0 * x^{m}\right)=0 * x^{-m}$. In effect, $x^{m}=0 * x^{-m}$ for any $m \in \mathbb{Z}$. Furthermore, in [7], we have $x^{m} * x^{n}=x^{m-n},\left(x^{m}\right)^{n}=x^{m n}$ for all $m, n \in \mathbb{Z}$, and $\langle x\rangle_{B}=\left\{x^{n}: n \in \mathbb{Z}\right\}$. If there exists a positive integer $n$ such that $x^{n}=0$, then the smallest such positive integer is denoted by $|x|_{B}$. If no such positive integer $n$ exists, then we say that $|x|_{B}$ is infinite. If $A \subseteq X$, then we denote $|A|_{B}$ as the cardinality of $A$.

Let $H$ and $K$ be subalgebras of $X$. In [4], we define the subset $H K$ of $X$ to be the set $H K=\{x \in X: x=h *(0 * k)$ for some $h \in H, k \in K\}$. Clearly, we have $H \subseteq H K$, $H \subseteq K H, K \subseteq H K$, and $K \subseteq K H$. Moreover, if $H \subseteq K$, then $H K=K H=K$. Also, $H K$ is a subalgebra of $X$ if and only if $H K=K H$ if and only if $H K=\langle H \cup K\rangle_{B}$. A B-algebra $X$ is called a cyclic B-algebra [7] if there exists $x \in X$ such that $X=\langle x\rangle_{B}$. Every cyclic B-algebra is commutative, but the converse need not be true. In [5], if $X=\langle x\rangle_{B}$ is a cyclic B-algebra with $|X|_{B}=m>1$ and if $H$ is a nontrivial subalgebra of $X$, then $H=\left\langle x^{k}\right\rangle_{B}$ for some integer $k>1$ such that $k$ divides $m$ and $|H|_{B}$ divides $m$. Furthermore, for every positive divisor $d$ of $m$, there exists a unique subalgebra $H$ of $X$ with $|H|_{B}=d$.

Let $H$ be a subalgebra of $X$ and $x \in X$. Let $x H=\{x *(0 * h): h \in H\}$ and $H x=\{h *(0 * x): h \in H\}$, called the left and right B-cosets of $H$ in $X$, respectively. If $X$ is commutative, then $x H=H x$ for all $x \in X$. Observe that $0 H=H=H 0$ and $x=x *(0 * 0) \in x H$ and $x=0 *(0 * x) \in H x$. It is easy to see that $x H=H$ if and only if $x \in H$.

Theorem 2.3. [2] Let $H$ be a subalgebra of $X$ and $a, b \in X$. Then
i. $a H=b H$ if and only if $(0 * b) *(0 * a) \in H$
ii. $H a=H b$ if and only if $a * b \in H$.

In [2], if $H$ is a subalgebra of $X$, then $\{x H: x \in X\}$ forms a partition of $X$ and there is a one-one correspondence of the set of all left B-cosets of $H$ in $X$ onto the set of all right B-cosets of $H$ in $X$. Thus, we define the number of distinct left (or right) B-cosets, written [ $X: H]_{B}$, of $H$ in $X$ as the index of $H$ in $X$. If $X$ is finite, then clearly $[X: H]_{B}$ is finite.

Theorem 2.4. [2] (Lagrange's Theorem for B-algebras) Let $H$ be a subalgebra of a finite $B$-algebra $X$. Then $|X|_{B}=[X: H]_{B}|H|_{B}$.

Corollary 2.5. [2] Let $|X|_{B}=p$, where $p$ is prime. Then $X$ is cyclic.
Theorem 2.6. [2] If $H, K$ are finite subalgebras of $X$, then $|H K|_{B}=\frac{|H|_{B}|K|_{B}}{|H \cap K|_{B}}$.
3 Some Implications of Lagrange's Theorem for B-algebras We now prove some results where Lagrange's Theorem plays a role.

Proposition 3.1. Let $X$ be a noncyclic B-algebra with $|X|_{B}=p^{2}$, where $p$ is prime. Then $|x|_{B}=p$ for every nonzero $x \in X$.

Proof. Let $x \in X$ and $x \neq 0$. By Lagrange's Theorem, $|x|_{B}$ divides $|X|_{B}=p^{2}$. Hence, $|x|_{B}$ is equal to $1, p$, or $p^{2}$. If $|x|_{B}=p^{2}$, then $\langle x\rangle_{B}=X$ and so $X$ is cyclic, a contradiction. Since $x \neq 0,|x|_{B} \neq 1$. Thus, $|x|_{B}=p$.
Proposition 3.2. If $X$ is a B-algebra with prime order, then $X$ has only the trivial subalgebras.

Proof. Suppose that $|X|_{B}=p$, where $p$ is prime. Let $H$ be a subalgebra of $X$. By Lagrange's Theorem, $|H|_{B}$ is 1 or $p$. Thus, $H=\{0\}$ or $H=X$.

Proposition 3.3. Let $|X|_{B}=p^{n}$, where $p$ is prime and $n \geq 1$. Then $X$ contains an element of order $p$.

Proof. Let $x \in X$ and $x \neq 0$. Then $H=\langle x\rangle_{B}$ is a cyclic subalgebra of $X$. By Lagrange's Theorem, $|H|_{B}$ divides $|X|_{B}=p^{n}$. Hence, $|H|_{B}=p^{m}$ for some $m \in \mathbb{Z}, 0<m \leq n$. It follows that for every divisor $d$ of $p^{m}$, there exists a subalgebra of order $d$. In particular, for $p$, there exists a subalgebra $K$ of $H$ such that $|K|_{B}=p$. By Corollary $2.5, K$ is cyclic and so there exists $y \in K$ such that $K=\langle y\rangle_{B}$ and $y$ is of order $p$. Hence, $X$ contains an element of order $p$.

Proposition 3.4. Let $X$ be a finite commutative $B$-algebra such that $X$ contains two distinct elements of order 2. Then $|X|_{B}$ is a multiple of 4.

Proof. Let $x$ and $y$ be two distinct elements of order 2. Let $H=\{0, x\}$ and $K=\{0, y\}$. Now, $H$ and $K$ are subalgebras of $X$. Since $X$ is commutative, $H K=\{0, x, y, x *(0 * y)\}$ is a subalgebra of $X$ of order 4. By Lagrange's Theorem, $|H K|_{B}=4$ divides $|X|_{B}$. Thus, $|X|_{B}$ is a multiple of 4 .

The above result need not be true if $X$ is not commutative. For instance, consider the B-algebra $X=\{0,1,2,3,4,5\}$ in Example 2.2. Note that $X$ is not commutative. Now, 3 and 4 are elements of $X$ with $|3|_{B}=2$ and $|4|_{B}=2$. However, 4 does not divide $|X|_{B}=6$.

Proposition 3.5. Let $X$ be a B-algebra with $|X|_{B}=p q$, where $p$ and $q$ are prime numbers. Then every proper subalgebra of $X$ is cyclic.

Proof. Let $H$ be a proper subalgebra of $X$. By Lagrange's Theorem, $|H|_{B}$ is $1, p, q$, or $p q$. Since $H$ is proper, $|H|_{B}$ is $p$ or $q$. By Corollary $2.5, H$ is cyclic.

Proposition 3.6. Let $H$ and $K$ be subalgebras of a finite $B$-algebra $X$ such that $|H|_{B}>\sqrt{|X|_{B}}$ and $|K|_{B}>\sqrt{|X|_{B}}$. Then $|H \cap K|_{B}>1$.

Proof. Suppose that $H$ and $K$ are subalgebras of a finite B-algebra $X$ such that $|H|_{B}>\sqrt{|X|_{B}}$ and $|K|_{B}>\sqrt{|X|_{B}}$. By Theorem 2.6, $|H \cap K|_{B}=\frac{|H|_{B}|K|_{B}}{|H K|_{B}}$. Since $|H|_{B}>\sqrt{|X|_{B}}$ and $|K|_{B}>\sqrt{|X|_{B}}$, it follows that $|H|_{B}|K|_{B}>|X|_{B}$. Since $|H K|_{B} \leq|X|_{B}$, it follows that $\frac{|X|_{B}}{|H K|_{B}} \geq 1$. Therefore, $|H \cap K|_{B}=\frac{|H|_{B}|K|_{B}}{|H K|_{B}}>\frac{|X|_{B}}{|H K|_{B}} \geq 1$.
Proposition 3.7. Let $|X|_{B}=p q$, where $p$ and $q$ are distinct primes with $p>q$. Then $X$ has at most one subalgebra of order $p$.
Proof. Suppose that $H$ and $K$ are subalgebras with $|H|_{B}=p=|K|_{B}$. Then $|H|_{B}>\sqrt{|X|_{B}}$ and $|K|_{B}>\sqrt{|X|_{B}}$. By Proposition 3.6, $|H \cap K|_{B}>1$. Thus, $|H \cap K|_{B}=p$ and so $H=K$.

4 Cauchy's Theorem for B-algebras This section establishes the Cauchy's Theorem for B-algebras and it also provides some implications of this theorem. We start with a simple observation given in the following lemma.

Lemma 4.1. Let $a \in X$. Then $a \in Z(X)$ if and only if $[X: C(a)]_{B}=1$ if and only if $C(a)=X$.

Let $a \in X$. An element $b \in X$ is said to be a conjugate of $a$ in $X$ if there exists $c \in X$ such that $b=c *(c * a)$. Let $R=\{(a, b) \in X \times X: b$ is a conjugate of $a\}$.

Theorem 4.2. Let $a \in X$. Then the relation $R$ on $X$ is an equivalence relation.
Proof. Since $a=0 *(0 * a)$, $a$ is conjugate to $a$. Thus, $R$ is reflexive. Let $(a, b) \in R$. Then there exists $c \in X$ such that $b=c *(c * a)$. Multiplying both sides by $0 * c$ twice, we have $(0 * c) *((0 * c) * b)=(0 * c) *[(0 * c) *(c *(c * a))]$. By (P2), (P3), (I), and (P1), we obtain

$$
\begin{aligned}
(0 * c) *((0 * c) * b) & =(0 * c) *[(0 * c) *(c *(c * a))] \\
& =(0 * c) *[((0 * c) *(0 *(c * a))) * c] \\
& =(0 * c) *[((0 * c) *(a * c)) * c] \\
& =(0 * c) *[(0 * a) * c)] \\
& =((0 * c) *(0 * c)) *(0 * a) \\
& =0 *(0 * a) \\
& =a
\end{aligned}
$$

Hence, $a$ is conjugate to $b$. Thus, $R$ is symmetric. Let $(a, b),(b, c) \in R$. Then there exist $u, v \in X$ such that $b=u *(u * a)$ and $c=v *(v * b)$. Now, by (P2) and (P3), we obtain

$$
\begin{aligned}
c & =v *(v * b) \\
& =v *[v *(u *(u * a))] \\
& =v *[(v *(0 *(u * a))) * u] \\
& =v *[(v *(a * u)) * u] \\
& =(v *(0 * u)) *(v *(a * u)) \\
& =(v *(0 * u) *[(v *(0 * u)) * a]
\end{aligned}
$$

Hence, $(a, c) \in R$ and so $R$ is transitive. Therefore, $R$ is an equivalence relation on $X$.
The equivalence relation $R$ in Theorem 4.2 is called conjugacy on $X$. The equivalence class of $a \in X$, denoted by $[a]_{c}$, of the relation $R$ is called the conjugacy class of $a$ in $X$.

Example 4.3. Consider the B-algebra $X=\{0,1,2,3,4,5\}$ in Example 2.2. Then there are three distinct conjugacy classes in $X$, namely, $[0]_{c}=\{0\},[1]_{c}=[2]_{c}=\{1,2\},[3]_{c}=[4]_{c}=$ $[5]_{c}=\{3,4,5\}$.

Remark 4.4. Let $a \in X$. Then $a \in Z(X)$ if and only if $[a]_{c}=\{a\}$.
The following theorem shows that the number of conjugates of $a$ is equal to the index of $C(a)$ in $X$.

Theorem 4.5. Let $a \in X$. Then $\left|[a]_{c}\right|_{B}=[X: C(a)]_{B}$.

Proof. Let $a \in X$. Let $\mathcal{L}$ denote the set of all distinct left B-cosets of $C(a)$ in $X$. Then $|\mathcal{L}|_{B}=[X: C(a)]_{B}$. By definition, $b *(b * a) \in[a]_{c}$ for all $b \in X$. Define $f: \mathcal{L} \rightarrow[a]_{c}$ by $f(b C(a))=b *(b * a)$. Suppose that $f(b C(a))=f(c C(a))$. Then by (P2), (P3), (P5), (I), (III), and Theorem 2.3(i), we have

$$
\begin{aligned}
f(b C(a))=f(c C(a)) & \Rightarrow b *(b * a)=c *(c * a) \\
& \Rightarrow 0 *(b *(b * a))=0 *(c *(c * a)) \\
& \Rightarrow(b * a) * b=(c * a) * c \\
& \Rightarrow(0 * c) *((b * a) * b)=(0 * c) *((c * a) * c) \\
& \Rightarrow(0 * c) *((b * a) * b)=((0 * c) *(0 * c)) *(c * a) \\
& \Rightarrow(0 * c) *((b * a) * b)=0 *(c * a) \\
& \Rightarrow(0 * c) *((b * a) * b)=a * c \\
& \Rightarrow[(0 * c) *((b * a) * b)] *(0 * b)=(a * c) *(0 * b) \\
& \Rightarrow[((0 * c) *(0 * b)) *(b * a)] *(0 * b)=a *((0 * b) *(0 * c)) \\
& \Rightarrow((0 * c) *(0 * b)) *[(0 * b) *(0 *(b * a))]=a *[0 *((0 * c) *(0 * b))] \\
& \Rightarrow((0 * c) *(0 * b)) *((0 * b) *(a * b))=a *[0 *((0 * c) *(0 * b))] \\
& \Rightarrow((0 * c) *(0 * b)) *(0 * a)=a *[0 *((0 * c) *(0 * b))] \\
& \Rightarrow(0 * c) *(0 * b) \in C(a) \\
& \Rightarrow b C(a)=c C(a) .
\end{aligned}
$$

Therefore, $f$ is a one-one function. Let $y \in[a]_{c}$. Then there exists $x \in X$ such that $y=x *(x * a)=f(x C(a))$. Hence, $f$ is onto. Therefore, $f$ is a one-one function from $\mathcal{L}$ onto $[a]_{c}$. Consequently, $\left|[a]_{c}\right|_{B}=|\mathcal{L}|_{B}=[X: C(a)]_{B}$.

Corollary 4.6. Let $X$ be a finite B-algebra. Then $|X|_{B}=\sum_{a}[X: C(a)]_{B}$, where the summation is over a complete set of distinct conjugacy class representatives.
Proof. By Theorem 4.2, $X=\bigcup_{a}[a]_{c}$, where the union runs over a complete set of distinct conjugacy class representatives. Since the distinct conjugacy classes are mutually disjoint, we have $|X|_{B}=\left|\bigcup_{a}[a]_{c}\right|_{B}=\sum_{a}\left|[a]_{c}\right|_{B}$. By Theorem 4.5, it follows that $|X|_{B}=\sum_{a}[X: C(a)]_{B}$, where the summation is over a complete set of distinct conjugacy class representatives.

Consider the B-algebra $X=\{0,1,2,3,4,5\}$ in Example 2.2. Then $|X|_{B}=6=1+2+3=$ $\left|[0]_{C}\right|_{B}+\left|[1]_{c}\right|_{B}+\left|[3]_{c}\right|_{B}=\sum_{a}\left|[a]_{c}\right|_{B}=\sum_{a}[X: C(a)]_{B}$.
Corollary 4.7. If $X$ is a finite B-algebra, then $|X|_{B}=|Z(X)|_{B}+\sum_{a \notin Z(X)}[X: C(a)]_{B}$, where the summation runs over a complete set of distinct conjugacy class representatives, which do not belong to $Z(X)$.
Proof. By Corollary 4.6, $|X|_{B}=\sum_{a}[X: C(a)]_{B}$, where the summation is over a complete set of distinct conjugacy class representatives. Thus, we have $|X|_{B}=\sum_{a \in Z(X)}[X: C(a)]_{B}+$
$\sum_{a \notin Z(X)}[X: C(a)]_{B}$. By Lemma 4.1, we have $\sum_{a \in Z(X)}[X: C(a)]_{B}=|Z(X)|_{B}$. Hence, $|X|_{B}=|Z(X)|_{B}+\sum_{a \notin Z(X)}[X: C(a)]_{B}$, where the summation runs over a complete set of distinct conjugacy class representatives which do not belong to $Z(X)$.

Example 4.8. Consider the B-algebra $X=\{0,1,2,3,4,5\}$ in Example 2.2. Then $Z(X)=$ $\{0\}$. Hence, $|Z(X)|_{B}+\sum_{a \notin Z(X)}[X: C(a)]_{B}=1+\left|[1]_{c}\right|_{B}+\left|[3]_{c}\right|_{B}=1+2+3=6=|X|_{B}$.

We now prove a partial converse of Lagrange's Theorem.
Lemma 4.9. If $X$ is a finite commutative B-algebra with $|X|_{B}=n$ such that $n$ is divisible by a prime $p$, then $X$ contains an element of order $p$ and hence a subalgebra of order $p$.

Proof. We proceed by induction on the order of $X$. If $|X|_{B}=p$ where $p$ is prime, then every element of $X$ (except 0 ) has order $p$. Thus, in particular, the lemma is true when $|X|_{B}=2$. Suppose that the lemma is true for all B-algebras of order $r$, where $2 \leq r<n$. Suppose that $X$ is a B-algebra of order $n$. Let $a \in X$ with $a \neq 0$ and let $|a|_{B}=m$. Then either $p \mid m$ or $p \nmid m$. If $p \mid m$, then $m=p k$ for some $k \in \mathbb{Z}^{+}$. In this case, $\left(a^{k}\right)^{p}=a^{m}=0$. Hence, $a^{k} \neq 0$ and $\left|a^{k}\right|_{B}=p$. Suppose $p \nmid m$. Since $X$ is commutative, the cyclic subalgebra $H=\langle a\rangle_{B}$ of $X$ is a normal subalgebra of $X$. By Lagrange's Theorem, $|X|_{B}=m[X: H]_{B}$. Since $p \nmid m$, we have $p\left|[X: H]_{B}=|X / H|_{B}\right.$. Since $| X /\left.H\right|_{B}<n$, there exists $b H \in X / H$ s.t. $|b H|_{B}=p$. Now, $b^{p} H=(b H)^{p}=H$. Hence, $b^{p} \in H$. Thus, $\left(b^{m}\right)^{p}=\left(b^{p}\right)^{m}=0$ and so either $b^{m}=0$ or $\left|b^{m}\right|_{B}=p$. If $b^{m}=0$, then $(b H)^{m}=H$ which implies $p \mid m$, a contradiction. Therefore, $\left|b^{m}\right|_{B}=p$ and so $b^{m}$ is the desired element of $X$.

Theorem 4.10. (Cauchy's Theorem for $B$-algebras) Let $X$ be a finite $B$-algebra with $|X|_{B}=n$ such that $n$ is divisible by a prime $p$. Then $X$ contains an element of order $p$ and hence a subalgebra of order $p$.

Proof. We proceed by induction on the order of $X$. If $n=2$, then $X$ is commutative and the result follows from Lemma 4.9. Suppose that the theorem is true for all B-algebras of order $m$ s.t. $2 \leq m<n$. By Corollary 4.7, $|X|_{B}=|Z(X)|_{B}+\sum_{a \notin Z(x)}[X: C(a)]$. If $X=Z(X)$, then $X$ is commutative and the result follows from Lemma 4.9. If $X \neq Z(X)$, then there exists $a \in X$ s.t. $a \notin Z(X)$. Then $X \neq C(a)$ and so $[X: C(a)]_{B}>1$. By Lagrange's Theorem, $|X|_{B}=[X: C(a)]_{B}|C(a)|_{B}>|C(a)|_{B}$. If $p\left||C(a)|_{B}\right.$, then $C(a)$ has an element of order $p$ and so $X$ has an element of order $p$. If $p \nmid|C(a)|_{B}$ for all $a \notin Z(X)$ , then $p \mid[X: C(a)]_{B}$ for all $a \notin Z(X)$. Since $p$ divides each term of the summation and also divides $|X|_{B}$, we have $p||Z(X)|$. By Lemma 4.9, $X$ contains an element of order $p$ and hence a subalgebra of order $p$

The following theorem proves that the converse of Lagrange's Theorem for B-algebras hold for finite commutative B-algebras.

Theorem 4.11. Let $X$ be a finite commutative B-algebra with $|X|_{B}=n$. If $m \in \mathbb{Z}^{+}$such that $m \mid n$, then $X$ has a subalgebra of order $m$.

Proof. If $m=1$, then $\{0\}$ is the required subalgebra of order $m$. If $n=1$, then $m=n=1$ and the result follows easily. Assume that $m>1$ and $n>1$. We proceed by induction on $n$. If $n=2$, then $m=2$ and $X$ is the required subalgebra of order $m$. Suppose that the theorem is true for all finite commutative B-algebras of order $k$ s.t. $2 \leq k<n$. Let
$p$ be a prime integer s.t. $p \mid m$. Then there exists $m_{1} \in \mathbb{Z}^{+}$s.t. $m=p m_{1}$. By Cauchy's Theorem, $X$ has a subalgebra $H$ of order $p$. Since $X$ is commutative, $H$ is normal and $X / H$ is a B-algebra. Now, $1 \leq|X / H|_{B}=\frac{|X|_{B}}{|H|_{B}}<|X|_{B}$ and $|X / H|_{B}=\frac{n}{p}$. Now, $n=m m_{2}$ for some $m_{2} \in \mathbb{Z}^{+}$. Thus, $|X / H|_{B}=\frac{p m_{1} m_{2}}{p}=m_{1} m_{2}$ and so $m_{1}$ divides $|X / H|_{B}$. Hence, $X / H$ has a subalgebra $K / H$ s.t. $|K / H|_{B}=m_{1}$, where $K$ is a subalgebra of $X$. Now, $|K|_{B}=|K / H|_{B}|H|_{B}=m_{1} p=m$. Hence, $K$ is a subalgebra of order $m$.

As a consequence of Cauchy's Theorem, we now introduce the concept of $\mathrm{B}_{p}$-algebras.
Definition 4.12. Let $p$ be a prime number. A B-algebra $X$ is called a $B_{p}$-algebra if the order of each element of $X$ is a power of $p$. A subalgebra $H$ of a B-algebra $X$ is called $B_{p}$-subalgebra if $H$ is a $\mathrm{B}_{p}$-algebra.

The B-algebra in Example 2.1 is $\mathrm{B}_{3}$-algebra. We now prove some results where Cauchy's Theorem plays a role. The following theorem provides a necessary and sufficient condition for a finite B -algebra to be a $\mathrm{B}_{p}$-algebra.

Theorem 4.13. Let $X$ be a nontrivial B-algebra. Then $X$ is a finite $B_{p}$-algebra if and only if $|X|_{B}=p^{k}$ for some $k \in \mathbb{Z}^{+}$.

Proof. Suppose that $X$ is a finite $\mathrm{B}_{p}$-algebra. If $q \|\left. X\right|_{B}$ for some prime $q \neq p$, then by Cauchy's Theorem, $X$ has an element of order $q$, a contradiction. Thus, $p$ is the only prime divisor of $|X|_{B}$, that is, $|X|_{B}=p^{k}$ for some $k \in \mathbb{Z}^{+}$. Conversely, suppose that $|X|_{B}=p^{k}$ for some $k \in \mathbb{Z}^{+}$. Then by Lagrange's Theorem, the order of each element of $X$ is a power of $p$. Therefore, $X$ is a finite $\mathrm{B}_{p}$-algebra.

The following theorem shows that the center of a $B_{p}$-algebra is nontrivial.
Theorem 4.14. If $X$ is a finite $B_{p}$-algebra with $|X|_{B}>1$, then $|Z(G)|_{B}>1$.
Proof. Suppose that $X$ is a finite $B_{p}$-algebra with $|X|_{B}>1$. If $X=Z(X)$, then $|Z(X)|_{B}=$ $|X|_{B}>1$. Suppose that $Z(X) \subset X$ and consider $a \in X$ such that $a \notin Z(X)$. Then $C(a)$ is a proper subalgebra of a $\mathrm{B}_{p}$-algebra $X$. By Theorem 4.13, $p \|\left. X\right|_{B}$. It follows that $p \mid[X: C(a)]_{B}$ for all $a \notin Z(X)$. Thus, $p$ divides $\sum_{a \notin Z(X)}[X: C(a)]_{B}$. By Corollary 4.7, $|X|_{B}=|Z(X)|_{B}+\sum_{a \notin Z(X)}[X: C(a)]_{B}$. Since $p \|\left. X\right|_{B}$ and $p \mid \sum_{a \notin Z(X)}[X: C(a)]_{B}$, it follows that $p \|\left. Z(X)\right|_{B}$. Therefore, $|Z(X)|_{B}>1$.

Corollary 4.15. If $|X|_{B}=p^{2}$, where $p$ is prime, then $X$ is commutative.
Proof. Suppose that $|X|_{B}=p^{2}$, where $p$ is prime. By Theorem 4.14, $|Z(X)|_{B}>1$. Since $Z(X)$ is a subalgebra, $|Z(X)|_{B}$ divides $p^{2}$ by Lagrange's Theorem. Hence, $|Z(X)|_{B}$ is $p$ or $p^{2}$. If $|Z(X)|_{B}=p$. Then $Z(X) \neq X$ and so there exists $a \in X$ such that $a \notin Z(X)$. In [6], $C(a)$ is a subalgebra of $X$ with $a \in C(a)$. Hence, $Z(X) \subset C(a)$. This implies that $|C(a)|_{B}=p^{2}$. Thus, $X=C(a)$ and so $a \in Z(X)$, a contradiction. Therefore, $|Z(X)|_{B}=p^{2}$ and so $X=Z(X)$. Consequently, $X$ is commutative.

Proposition 4.16. Let $H$ and $K$ be subalgebras of a commutative $B$-algebra $X$. If $|H|_{B}=$ $m$ and $|K|_{B}=n$, then $X$ has a subalgebra of $\operatorname{order} \operatorname{lcm}(m, n)$.

Proof. Let $H$ and $K$ be subalgebras of a commutative B-algebra $X$ with $|H|_{B}=m$ and $|K|_{B}=n$. Since $H K=K H, H K$ is a subalgebra of $X$. Since $H$ and $K$ are finite, $H$ and $K$ are subalgebras of a finite B-algebra $H K$. By Lagrange's Theorem, $m \|\left. H K\right|_{B}$ and $n \|\left. H K\right|_{B}$. Hence, $l c m(m, n) \|\left. H K\right|_{B}$. By Theorem 4.11, $H K$ has a subalgebra of order $\operatorname{lcm}(m, n)$ and so $X$ has a subalgebra of order $\operatorname{lcm}(m, n)$.

The version of Lagrange's Theorem for B-algebras in [2] is analogue to the Lagrange's Theorem for groups, and the version of Cauchy's Theorem for B-algebras in this paper is analogue to the Cauchy's Theorem for groups. It is then natural to seek an analogue results to the Sylow Theorems for groups.

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# TRANSFORMS ON OPERATOR MONOTONE FUNCTIONS 

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Abstract. Let $f$ be an operator monotone function on $[0, \infty)$ with $f(t) \geq 0$ and $f(1)=1$. If $f(t)$ is neither the constant function 1 nor the identity function $t$, then

$$
h(t)=\frac{(t-a)(t-b)}{(f(t)-f(a))\left(f^{\sharp}(t)-f^{\sharp}(b)\right)} \quad t \geq 0
$$

is also operator monotone on $[0, \infty)$, where $a, b \geq 0$ and

$$
f^{\sharp}(t)=\frac{t}{f(t)} \quad t \geq 0 .
$$

Moreover, we show some extensions of this statement.
1 Introduction and History We call a real continuous function $f(t)$ on an interval $I$ operator monotone on $I$ (in short, $f \in \mathbb{P}(I)$ ), if $A \leq B$ implies $f(A) \leq f(B)$ for any self-adjoint matrices $A, B$ with their spectrum containd in $I$. In this paper, we consider only the case $I=[0, \infty)$ or $I=(0, \infty)$. We denote $f \in \mathbb{P}_{+}(I)$ if $f \in \mathbb{P}(I)$ satisfies $f(t) \geq 0$ for any $t \in I$.

Let $\mathbb{H}_{+}$be the upper-half plain of $\mathbb{C}$, that is,

$$
\mathbb{H}_{+}=\{z \in \mathbb{C} \mid \operatorname{Im} z>0\}=\{z \in \mathbb{C}| | z \mid>0,0<\arg z<\pi\},
$$

where $\operatorname{Im} z$ (resp. $\arg z)$ means the imaginary part (resp. the argument) of $z$. As Loewner's theorem, it is known that $f$ is operator monotone on $I$ if and only if $f$ has an analytic continuation to $\mathbb{H}_{+}$that maps $\mathbb{H}_{+}$into its closure $\overline{\mathbb{H}}_{+}$ and also has an analytic continuation to the lower half-plane $\mathbb{H}_{-}$, obtained by the reflection across $I$. (see [1],[3],[5] ).
D. Petz [11] proved that an operator monotone function $f:[0, \infty) \longrightarrow$ $[0, \infty)$ satisfying the functional equation

$$
f(t)=t f\left(t^{-1}\right) \quad t \geq 0
$$

[^4]is related to a Morozova-Chentsov function [9] which gives a monotone metric on the manifold of $n \times n$ density matrices. In the work [12], the concrete functions (Petz-Hasegawa's functions)
$$
f_{a}(t)=a(1-a) \frac{(t-1)^{2}}{\left(t^{a}-1\right)\left(t^{1-a}-1\right)} \quad(-1 \leq a \leq 2)
$$
appeared and their operator monotonicity was proved (see also [2]). V.E.S. Szabo introduced an interesting idea for checking their operator monotonicity in [13], but his idea was something strange. We use a similar idea in our argument. M. Uchiyama [14] proved the operator monotonicity of the following extended functions:
$$
\frac{(t-a)(t-b)}{\left(t^{p}-a^{p}\right)\left(t^{1-p}-b^{1-t}\right)} \in \mathbb{P}_{+}[0 . \infty)
$$
for $0<p<1$ and $a, b>0$. The main result of this paper is as follows:
Theorem 1. Let $a$ and $b$ be non-negative real numbers. If $f \in \mathbb{P}_{+}[0, \infty)$ and both $f$ and $f^{\sharp}$ are not constant, then
$$
h(t)=\frac{(t-a)(t-b)}{(f(t)-f(a))\left(f^{\sharp}(t)-f^{\sharp}(b)\right)} \in \mathbb{P}_{+}[0 . \infty),
$$
where
$$
f^{\sharp}(t)=\frac{t}{f(t)} \quad t \geq 0
$$

The proof of this statement was given in [7] by the author and his student, M. Kawasaki. We also made its revised version as [8]. But these manuscripts have been unpublished. F. Hansen [4] has inspired by [8] and given the different proof of the above statement in the case $a=b=1$. The proof in this paper was based on the theory of Complex Analysis and we have considered this method useful. The same method was introduced for the proof of the operator monotonicity of Petz-Hasegawa's functions in [6]. Also the autor and S. Wada have extended the method and succeeded to justify Szabo's result in a sense [10]. The last section we will prove main theorem and give some applications.

2 Main result For $f \in \mathbb{P}[0, \infty)$, we have the following integral representation:

$$
f(z)=f(0)+\beta z+\int_{0}^{\infty} \frac{\lambda z}{z+\lambda} d w(\lambda)
$$

where $\beta \geq 0$ and

$$
\int_{0}^{\infty} \frac{\lambda}{1+\lambda} d w(\lambda)<\infty
$$

(see [1]). When $f(0) \geq 0$ (i.e., $f \in \mathbb{P}_{+}[0, \infty)$ ), it holds that $0<\arg f(z) \leq$ $\arg z$ whenever $0<\arg z<\pi$.

For any $f \in \mathbb{P}_{+}[0, \infty)(f \neq 0)$, we define $f^{\sharp}$ as follows:

$$
f^{\sharp}(t)=\frac{t}{f(t)} \quad t \in[0, \infty) .
$$

Then it is well-known $f^{\sharp} \in \mathbb{P}_{+}[0, \infty)$.

Proposition 2. Let $f$ be an operator monotone function on $(0, \infty)$ and $a$ be a positive real number.
(1) When $f(t)$ is not constant, we have

$$
g_{1}(t)=\frac{t-a}{f(t)-f(a)} \in \mathbb{P}_{+}[0, \infty)
$$

(2) When $f(t) \geq 0$ for $t \geq 0$, we have

$$
g_{2}(t)=\frac{f(t)(t-a)}{t f(t)-a f(a)} \in \mathbb{P}_{+}[0, \infty)
$$

Proof. (1) It has proved in [14]. We state the outline of the proof. For $f \in$ $\mathbb{P}_{+}(0, \infty)$, we have

$$
f(z)=\alpha+\beta z+\int_{0}^{\infty}\left(-\frac{1}{x+z}+\frac{x}{x^{2}+1}\right) d \nu(x) \quad(\alpha \in \mathbb{R}, \beta \geq 0)
$$

for $z \in \mathbb{H}_{+} \cup(0, \infty)$ by Loewner's theorem ([1], [3], [6]). Since

$$
g_{1}(z)=\frac{z-a}{f(z)-f(a)}=\frac{1}{\beta+\int_{0}^{\infty} \frac{1}{(x+z)(x+a)} d \nu(x)},
$$

we have

$$
\begin{aligned}
\operatorname{Im} g_{1}(z) & =\frac{-1}{\left|g_{1}(z)\right|^{2}} \operatorname{Im}\left(\beta+\int_{0}^{\infty} \frac{1}{(x+z)(x+a)} d \nu(x)\right) \\
& =\frac{-1}{\left|g_{1}(z)\right|^{2}} \int_{0}^{\infty} \frac{1}{x+a} \operatorname{Im} \frac{1}{x+z} d \nu(x)>0
\end{aligned}
$$

for $z \in \mathbb{H}_{+}$. This implies $g_{1} \in \mathbb{P}_{+}[0, \infty)$.
(2) Since $g_{2}([0, \infty)) \subset[0, \infty)$, it suffices to show that $g_{2}\left(\mathbb{H}_{+}\right) \subset \mathbb{H}_{+}$. By the calculation

$$
\begin{aligned}
g_{2}(z) & =\frac{z f(z)-a f(a)+a f(a)-f(z) a}{z f(z)-a f(a)}=1-\frac{a(f(z)-f(a))}{z f(z)-a f(a)} \\
& =1-\frac{a}{\frac{z f(z)-a f(a)}{f(z)-f(a)}}=1-\frac{a}{z+f(a) g_{1}(z)}
\end{aligned}
$$

we have

$$
\operatorname{Im} g_{2}(z)=-\operatorname{Im} \frac{a}{z+f(a) g_{1}(z)}=\operatorname{Im} \frac{a\left(z+f(a) g_{1}(z)\right)}{\left|z+f(a) g_{1}(a)\right|^{2}}
$$

When $z \in \mathbb{H}_{+}, \operatorname{Im} g_{1}(z)>0$ by (1) and $\operatorname{Im} g_{2}(z)>0$. So $g_{2}(t)$ belongs to $\mathbb{P}_{+}[0, \infty)$.

Lemma 3. Let $a \geq 0$ and $0<p<1$. If $f$ is a non-constant operator monotone function on $(0, \infty)$, then we have

$$
f\left(t^{p}\right)-f\left(a^{p}\right) \neq 0 \text { and } t f\left(t^{p}\right)-a f\left(a^{p}\right) \neq 0
$$

for any $t \in(-\infty, 0)$.
Proof. When $f$ is operator monotone and not constant, we have $\operatorname{Im} f(z)>0$ for any $z \in \mathbb{H}_{+}$by the maximum principle for the harmonic function $\operatorname{Im} f$ on $\mathbb{H}_{+}$. For any $t=|t| e^{i \pi} \in(-\infty, 0)$, we have $t^{p} \in \mathbb{H}_{+}$and $\operatorname{Im} f\left(t^{p}\right)>0$. This implies

$$
\operatorname{Im} f\left(t^{p}\right) \neq \operatorname{Im} f\left(a^{p}\right)=0 \text { and } \operatorname{Im} f\left(t^{p}\right) \neq \operatorname{Im} \frac{a f\left(a^{p}\right)}{t}=0
$$

that is, $f\left(t^{p}\right)-f\left(a^{p}\right) \neq 0$ and $t f\left(t^{p}\right)-a f\left(a^{p}\right) \neq 0$.

Lemma 4. For any $z \in \mathbb{H}_{+}$and a positive integer $n(n \geq 2)$, we have

$$
\arg z<\arg (z-l)<\frac{\pi+(n-1) \arg z}{n} \quad \text { if } \quad 0<l \leq \frac{|z|}{n-1} .
$$

Proof. It is clear that $\arg z<\arg (z-l)$ for $z \in \mathbb{H}_{+}$and $l>0$.
It suffices to show that, for $z=e^{i \theta}(0<\theta<\pi)$,

$$
\arg (z-l)<\frac{\pi+(n-1) \theta}{n} \quad \text { if } 0<l \leq \frac{1}{n-1} .
$$

We set

$$
w=\frac{\sin \theta}{\sin \frac{\pi+(n-1) \theta}{n}} e^{i(\pi+(n-1) \theta) / n}
$$

Then we have $\operatorname{Im} z=\operatorname{Im} w$ and

$$
0<z-w=\cos \theta-\frac{\sin \theta}{\sin \frac{\pi+(n-1) \theta}{n}} \cos \frac{\pi+(n-1) \theta}{n}=\frac{\sin \frac{\pi-\theta}{n}}{\sin \frac{\pi+(n-1) \theta}{n}}
$$

By the estimation

$$
\begin{aligned}
& \inf \{z-w \mid 0<\theta<\pi\}=\inf \left\{\left.\frac{\sin \frac{\pi-\theta}{n}}{\sin \frac{\pi+(n-1) \theta}{n}} \right\rvert\, 0<\theta<\pi\right\} \\
= & \inf \left\{\left.\frac{\sin t}{\sin (\pi-(n-1) t)} \right\rvert\, 0<t<\frac{\pi}{n}\right\}=\inf \left\{\left.\frac{\sin t}{\sin (n-1) t} \right\rvert\, 0<t<\frac{\pi}{n}\right\}=\frac{1}{n-1},
\end{aligned}
$$

we can get the desired result.

Now we can prove the following theorem and remark that Theorem 1 follows from this statement because $f(t) f^{\sharp}(t) / t=1 \in \mathbb{P}_{+}[0, \infty)$ :

Theorem 5. Let $n$ be a positive integer, $a, b, b_{1}, \ldots, b_{n} \geq 0$ and $f, g, g_{1}, \ldots, g_{n}$ be non-constant, non-negative operator monotone functions on $[0, \infty)$.
(1) If $\frac{f(t) g(t)}{t}$ is operator monotone on $[0, \infty)$, then the function

$$
h(t)=\frac{(t-a)(t-b)}{(f(t)-f(a))(g(t)-g(b))}
$$

is operator monotone on $[0, \infty)$ for any $a, b \geq 0$.
(2) If $\frac{f(t)}{\prod_{i=1}^{n} g_{i}(t)}$ is operator monotone on $[0, \infty)$, then the function

$$
h(t)=\frac{(t-a)}{(f(t)-f(a))} \prod_{i=1}^{n} \frac{g_{i}(t)\left(t-b_{i}\right)}{t g_{i}(t)-b_{i} g_{i}\left(b_{i}\right)}
$$

is operator monotone on $[0, \infty)$ for any $a, b \geq 0$.
Proof. (1) By $f, g \in \mathbb{P}_{+}[0, \infty)$ and Proposition 2 (1),

$$
\frac{t-a}{f(t)-f(a)} \text { and } \frac{t-b}{g(t)-g(b)}
$$

are operator monotone on $[0, \infty)$. Therefore

$$
h(z)=\frac{(z-a)(z-b)}{(f(z)-f(a))(g(z)-g(b))}
$$

is holomorphic on $\mathbb{H}_{+}$, continuous on $\mathbb{H}_{+} \cup[0, \infty)$ and satisfies $h([0, \infty)) \subset$ $[0, \infty)$ and

$$
\arg h(z)=\arg \frac{z-a}{f(z)-f(a)}+\arg \frac{z-b}{g(z)-g(b)}>0 \text { for } z \in \mathbb{H}_{+} .
$$

We assume that $f(z)$ and $g(z)$ are continuous on the closure $\overline{\mathbb{H}_{+}}$of $\mathbb{H}_{+}$and

$$
f(t)-f(a) \neq 0 \text { and } g(t)-g(b) \neq 0 \text { for any } t \in(-\infty, 0)
$$

Then $h(z)$ is continuous on $\overline{\mathbb{H}_{+}}$.
In the case $z \in(-\infty, 0)$, i.e., $|z|>0$ and $\arg z=\pi$, we have

$$
\begin{aligned}
& \arg h(z) \\
= & \arg (z-a)-\arg (f(z)-f(a))+\arg (z-b)-\arg (g(z)-g(b)) \\
\leq & \pi-\arg f(z)+\pi-\arg g(z) \\
\leq & 2 \pi-\arg z=\pi \quad(\text { since } \arg f(z)+\arg g(z)-\arg z \geq 0) .
\end{aligned}
$$

So it holds $0 \leq \arg h(z) \leq \pi$.
In the case that $z \in \mathbb{H}_{+}$satisifying $|z|>\max \{a, b\}$, it holds that

$$
\arg (z-a), \arg (z-b)<\frac{\pi+\arg z}{2}
$$

by Lemma 4 (as $l=\max \{a, b\}$ and $n=2$ ). Since

$$
\begin{aligned}
\arg h(z) & =\arg (z-a)-\arg (f(z)-f(a))+\arg (z-b)-\arg (g(z)-g(b)) \\
& <\frac{\pi+\arg z}{2}-\arg f(z)+\frac{\pi+\arg z}{2}-\arg g(z) \\
& =\pi+\arg z-\arg f(z)-\arg g(z) \leq \pi
\end{aligned}
$$

we have $0<\arg h(z)<\pi$.
For $r>0$, we define $H(r)=\{z \in \mathbb{C}| | z \mid \leq r, \operatorname{Im} z \geq 0\}$. Whenever $r>l=\max \{a, b\}$, we can get

$$
0 \leq \arg h(z) \leq \pi
$$

on the boundary of $H(r)$. Since $h(z)$ is holomorphic on the interior $H(r)^{\circ}$ of $H(r)$ and continuous on $H(r), \operatorname{Im} h(z)$ is harmonic on $H(r)^{\circ}$ and continuous on
$H(r)$. Because $\operatorname{Im} h(z) \geq 0$ on the boundary of $H(r)$, we have $h(H(r)) \subset \overline{\mathbb{H}_{+}}$ by the minimum principle of harmonic functions. This implies

$$
h\left(\overline{\mathbb{H}_{+}}\right)=h\left(\bigcup_{r>l} H(r)\right) \subset \bigcup_{r>l} h(H(r)) \subset \overline{\mathbb{H}_{+}},
$$

and $h \in \mathbb{P}_{+}[0, \infty)$.
In general case, we define $f_{p}$ and $g_{p}$ as follows $(0<p<1)$ :

$$
f_{p}(t)=f\left(t^{p}\right) \text { and } g_{p}(t)=t^{1-p} g\left(t^{p}\right)
$$

It is clear $f_{p} \in \mathbb{P}_{+}[0, \infty)$. When $0<\arg z<\pi$, we have

$$
\begin{gathered}
\arg g_{p}(z) \geq(1-p) \arg z>0 \text { and } \\
\arg g_{p}(z)=(1-p) \arg z+\arg g\left(z^{p}\right) \leq(1-p) \arg z+p \arg z<\pi .
\end{gathered}
$$

So $g_{p} \in \mathbb{P}_{+}[0, \infty)$. For any $t \in(-\infty, 0)$, we have $f_{p}(t)-f_{p}(a) \neq 0$ by Lemma 3 and $g_{p}(t)-g_{p}(a) \neq 0$ because $t=|t| e^{\pi i}$ and $(1-p) \pi<\arg g_{p}(t) \leq \pi$. So we have

$$
h_{p}(z)=\frac{(z-a)(z-b)}{\left(f_{p}(z)-f_{p}(a)\right)\left(g_{p}(z)-g_{p}(b)\right)}
$$

is holomorphic on $\mathbb{H}_{+}$and continuous on $\overline{\mathbb{H}_{+}}$. Since $\frac{f(t) g(t)}{t} \in \mathbb{P}_{+}[0, \infty)$,

$$
\frac{f_{p}(t) g_{p}(t)}{t}=\frac{f\left(t^{p}\right) t^{1-p} g\left(t^{p}\right)}{t}=\frac{f\left(t^{p}\right) g\left(t^{p}\right)}{t^{p}}
$$

also belongs to $\mathbb{P}_{+}[0, \infty)$. By the above argument, we have $h_{p} \in \mathbb{P}_{+}[0, \infty)$.
Since

$$
\begin{aligned}
h_{p}(t) & =\frac{(t-a)(t-b)}{\left(f_{p}(t)-f_{p}(a)\right)\left(g_{p}(t)-g_{p}(b)\right)} \\
& =\frac{(t-a)(t-b)}{\left(f\left(t^{p}\right)-f\left(a^{p}\right)\right)\left(t^{1-p} g\left(t^{p}\right)-b^{1-p} g\left(b^{p}\right)\right)} \quad \text { for } t \geq 0,
\end{aligned}
$$

we have

$$
\lim _{p \rightarrow 1-0} h_{p}(t)=h(t) .
$$

So we can get the operator monotonicity of $h(t)$.
(2) We show this by the similar way as (1). By Proposition 2,

$$
\frac{t-a}{f(t)-f(a)} \text { and } \frac{g_{i}(t)\left(t-b_{i}\right)}{t g_{i}(t)-b_{i} g_{i}\left(b_{i}\right)} \quad(i=1,2, \ldots, n)
$$

are operator monotone on $[0, \infty)$. So we have that

$$
h(z)=\frac{z-a}{f(z)-f(a)} \prod_{i=1}^{n} \frac{g_{i}(z)\left(z-b_{i}\right)}{z g_{i}(z)-b_{i} g_{i}\left(b_{i}\right)}
$$

is holomorphic on $\mathbb{H}_{+}$, continuous on $\mathbb{H}_{+} \cup[0, \infty)$ and satisfies $h([0, \infty)) \subset$ $[0, \infty)$ and

$$
\arg h(z)=\arg \frac{z-a}{f(z)-f(a)}+\sum_{i=1}^{n} \arg \frac{g_{i}(z)\left(z-b_{i}\right)}{z g_{i}(z)-b_{i} g_{i}\left(b_{i}\right)}>0
$$

for $z \in \mathbb{H}_{+}$.
We assume that $f(z)$ and $g_{i}(z)(i=1,2, \ldots, n)$ are continuous on $\overline{\mathbb{H}_{+}}$and

$$
f(t)-f(a) \neq 0 \text { and } t g_{i}(t)-b_{i} g_{i}\left(b_{i}\right) \neq 0 \text { for any } t \in(-\infty, 0) .
$$

Then $h(z)$ is continuous on $\overline{\mathbb{H}_{+}}$.
In the case $z \in(-\infty, 0)$, i.e., $|z|>0$ and $\arg z=\pi$, we have

$$
\begin{aligned}
& \arg h(z) \\
&= \arg (z-a)+\sum_{i=1}^{n} \arg g_{i}(z)\left(z-b_{i}\right)-\arg (f(z)-f(a))-\sum_{i=1}^{n} \arg \left(z g_{i}(z)-b_{i} g_{i}\left(b_{i}\right)\right) \\
& \leq \pi+\sum_{i=1}^{n} \arg g_{i}(z)+n \pi-\arg f(z)-n \pi \quad\left(\text { since } \arg \left(z g_{i}(z)-b_{i} g_{i}\left(b_{i}\right)\right) \geq \pi\right) \\
& \leq \pi \quad \quad \quad\left(\text { since } \arg f(z)-\sum_{i=1}^{n} \arg g_{i}(z) \geq 0\right) .
\end{aligned}
$$

So it holds $0 \leq \arg h(z) \leq \pi$.
In the case $z \in \mathbb{H}_{+}$satisifying $|z|>n \max \left\{a, b_{1}, b_{2}, \ldots, b_{n}\right\}$, it holds that

$$
\arg (z-a), \arg \left(z-b_{i}\right)<\frac{\pi+n \arg z}{n+1} \quad(i=1,2, \ldots, n)
$$

by Lemma 4 . We may assume that there exists a number $k(1 \leq k \leq n)$ such that

$$
\arg \left(z g_{i}(z)\right) \leq \pi \quad(i \leq k), \quad \arg \left(z g_{i}(z)\right)>\pi \quad(i>k)
$$

Since

$$
\begin{aligned}
& \quad \arg h(z) \\
& =\arg (z-a)+\sum_{i=1}^{n} \arg \left(z-b_{i}\right)+\sum_{i=1}^{n} \arg g_{i}(z) \\
& \quad-\arg (f(z)-f(a))-\sum_{i=1}^{n} \arg \left(z g_{i}(z)-b_{i} g_{i}\left(b_{i}\right)\right) \\
& \leq \frac{\pi+n \arg z}{n+1} \times(n+1)+\sum_{i=1}^{n} \arg g_{i}(z) \\
& \quad \quad-\arg f(z)-\sum_{i=1}^{k} \arg z g_{i}(z)-(n-k) \pi \\
& = \\
& \pi+n \arg z+\sum_{i=k+1}^{n} \arg g_{i}(z)-\arg f(z)-k \arg z-(n-k) \pi \\
& \leq \pi+(n-k) \arg z-(n-k) \pi \leq \pi,
\end{aligned}
$$

we have $0 \leq \arg h(z) \leq \pi$.
This means that it holds

$$
0 \leq \arg h(z) \leq \pi
$$

if $z$ belongs to the boundary of $H(r)=\{z \in \mathbb{C}| | z \mid \leq r, \operatorname{Im} z \geq 0\}$ for a sufficiently large $r$. Using the same argument in (1), we can prove the operator monotonicity of $h$.

In general case, we define functions, for $p(0<p<1)$, as follows:

$$
f_{p}(t)=f\left(t^{p}\right), \quad g_{i, p}(t)=g_{i}\left(t^{p}\right) \quad(i=1,2, \ldots, n)
$$

Since $f, g_{i} \in \mathbb{P}_{+}[0, \infty)$,

$$
0<\arg f_{p}(z)<\pi, \quad 0<\arg z g_{i, p}(z)<2 \pi
$$

for $z \in \mathbb{H}_{+}$. This means that $f_{p}(z)$ and $g_{i, p}(z)$ are continuous on $\overline{\mathbb{H}_{+}}$and

$$
f_{p}(t)-f_{p}(a) \neq 0 \text { and } t g_{i, p}(t)-b_{i} g_{i, p}\left(b_{i}\right) \neq 0 \text { for any } t \in(-\infty, 0)
$$

by Lemma 3. Since

$$
\frac{f_{p}(t)}{\prod_{i=1}^{n} g_{i, p}(t)}=\frac{f\left(t^{p}\right)}{\prod_{i=1}^{n} g_{i}\left(t^{p}\right)} \quad(0<p<1)
$$

is operator monotone on $[0, \infty)$, we can get the operator monotonicity of

$$
\begin{aligned}
h_{p}(t) & =\frac{t-a}{f_{p}(t)-f_{p}(a)} \prod_{i=1}^{n} \frac{g_{i, p}(t)\left(t-b_{i}\right)}{t g_{i, p}(t)-b_{i} g_{i, p}\left(b_{i}\right)} \\
& =\frac{t-a}{f\left(t^{p}\right)-f\left(a^{p}\right)} \prod_{i=1}^{n} \frac{g_{i}\left(t^{p}\right)\left(t-b_{i}\right)}{t g_{i}\left(t^{p}\right)-b_{i} g_{i}\left(b_{i}^{p}\right)} .
\end{aligned}
$$

So we can see that

$$
h(t)=\lim _{p \rightarrow 1-0} h_{p}(t)
$$

is operator monotone on $[0, \infty)$.

Remark 6. Using Proposition 2 and Theorem 5, we can prove the operator monotonicity of the concrete functions in [12]. Since $t^{a}(0<a<1)$ and $\log t$ is operator monotone on $(0, \infty)$,

$$
\begin{aligned}
& f_{a}(t)=a(1-a) \frac{(t-1)^{2}}{\left(t^{a}-1\right)\left(t^{1-a}-1\right)} \quad(-1 \leq a \leq 2) \\
& = \begin{cases}a(a-1) \frac{t^{-a}(t-1)^{2}}{\left(t^{-a}-1\right)\left(t \cdot t^{-a}-1\right)} & -1 \leq a<0 \\
\frac{t-1}{\log t} & a=0,1 \\
a(1-a) \frac{(t-1)^{2}}{\left(t^{a}-1\right)\left(t^{1-a}-1\right)} & 0<a<1 \\
a(a-1) \frac{t^{-1}(t-1)^{2}}{\left(t^{a-1}-1\right)\left(t \cdot t^{a-1}-1\right)} & 1<a \leq 2\end{cases}
\end{aligned}
$$

becomes operator monotone.
We can also prove this remark and the first part of Example 10 using some formula stated in [10] as Theorem 1.2.

Corollary 7. Let $f \in \mathbb{P}_{+}(0, \infty)$ and both $f$ and $f^{\sharp}$ be not constant. For any $a>0$, we define

$$
h_{a}(t)=\frac{(t-a)\left(t-a^{-1}\right)}{(f(t)-f(a))\left(f^{\sharp}(t)-f^{\sharp}\left(a^{-1}\right)\right)} \quad t \in(0, \infty) .
$$

Then we have
(1) $h_{a}$ is operator monotone on $(0, \infty)$.
(2) $f(t)=t \cdot f\left(t^{-1}\right)$ implies $h_{a}(t)=t \cdot h_{a}\left(t^{-1}\right)$.
(3) $a=1$ and $f\left(t^{-1}\right)=f(t)^{-1}$ imply $h_{1}(t)=t \cdot h_{1}\left(t^{-1}\right)$.

Proof. We can directly prove (1) from Theorem 5. Because

$$
\begin{aligned}
t \cdot h_{a}\left(t^{-1}\right) & =\frac{t\left(t^{-1}-a\right)\left(t^{-1}-a^{-1}\right)}{\left(f\left(t^{-1}\right)-f(a)\right)\left(f^{\sharp}\left(t^{-1}\right)-f^{\sharp}\left(a^{-1}\right)\right)} \\
& =\frac{(t-a)\left(t-a^{-1}\right)}{t\left(f\left(t^{-1}\right)-f(a)\right)\left(f^{\sharp}\left(t^{-1}\right)-f^{\sharp}\left(a^{-1}\right)\right)},
\end{aligned}
$$

we can compute

$$
\begin{aligned}
& t\left(f\left(t^{-1}\right)-f(a)\right)\left(f^{\sharp}\left(t^{-1}\right)-f^{\sharp}\left(a^{-1}\right)\right)-(f(t)-f(a))\left(f^{\sharp}(t)-f^{\sharp}\left(a^{-1}\right)\right) \\
= & \left(f\left(t^{-1}\right)-f(a)\right)\left(1 / f\left(t^{-1}\right)-t / a f\left(a^{-1}\right)\right)-(f(t)-f(a))\left(t / f(t)-1 / a f\left(a^{-1}\right)\right) \\
= & 0
\end{aligned}
$$

if it holds $f(t)=t \cdot f\left(t^{-1}\right)$ or $a=1, f\left(t^{-1}\right)=f(t)^{-1}$. So we have (2) and (3).

The function $h$ is called symmetric if it satisfies the following condition:

$$
h(t)=t h\left(t^{-1}\right), \quad t \geq 0 .
$$

We can define a symmetric operator mean using a symmetric operator monotone function in the sense of Kubo-Ando ([5], [6]). Corollary 7 says that we can repeatedly construct a symmetric operator monotone function from a symmetric operator monotone function. We can give the following examples.

Example 8. If we choose $t^{p}(0<p<1)$ as $f(t)$ in Corollary 7(3),

$$
h(t)=\frac{(t-1)^{2}}{\left(t^{p}-1\right)\left(t^{1-p}-1\right)} .
$$

If we choose $t^{p}+t^{1-p}(0<p<1)$ as $f(t)$ in Corollary $7(2)$,

$$
\begin{aligned}
h(t) & =\frac{t-a}{t^{p}+t^{1-p}-a^{p}-a^{1-p}} \times \frac{t-a^{-1}}{\frac{1}{t^{p-1}+t^{-p}}-\frac{1}{a^{p}+a^{1-p}}} \quad(a>0) \\
& =\frac{\sqrt{t}(\cosh (\log t)-\cosh (\log a))}{\cosh (\log \sqrt{t})-\cosh \left(\log \sqrt{t}+\log \left(t^{p}+t^{1-p}\right)-\log \left(a^{p}+a^{1-p}\right)\right)} .
\end{aligned}
$$

These functions, $h \in \mathbb{P}_{+}[0, \infty)$, are symmetric.

3 Extension of Theorem 4 Let $m$ and $n$ be positive integers and $f_{1}, f_{2}, \ldots, f_{m}$, $g_{1}, g_{2}, \ldots, g_{n}$ be non-constant, non-negative operator monotone functions on $[0, \infty)$. We assume that the function

$$
F(t)=\frac{\prod_{i=1}^{m} f_{i}(t)}{t^{m-1} \prod_{j=1}^{n} g_{j}(t)}
$$

is operator monotone on $[0, \infty)$. For non-negative numbers $a_{1}, a_{2}, \ldots, a_{m}, b_{1}, b_{2}, \ldots, b_{n}$, we define the function $h(t)$ as follows:

$$
h(t)=\prod_{i=1}^{m} \frac{t-a_{i}}{f_{i}(t)-f_{i}\left(a_{i}\right)} \prod_{j=1}^{n} \frac{g_{j}(t)\left(t-b_{j}\right)}{t g_{j}(t)-b_{j} g_{j}\left(b_{j}\right)} \quad(t \geq 0)
$$

Then it follows from Proposition 2 that $h(z)$ is holomorphic on $\mathbb{H}_{+}, h([0, \infty)) \subset$ $[0, \infty)$ and $\arg h(z)>0$ for any $z \in \mathbb{H}_{+}$.

Theorem 9. In the above setting, we have the following:
(1) When $f_{i}$ and $g_{j}(1 \leq i \leq m, 1 \leq j \leq n)$ are continuous on $\overline{\mathbb{H}_{+}}$and

$$
f_{i}(t)-f_{i}\left(a_{i}\right) \neq 0, \quad t g_{j}(t)-b_{j} g_{j}\left(b_{j}\right) \neq 0, \quad t \in(-\infty, 0),
$$

$h(t)$ is operator monotone on $[0, \infty)$.
(2) When there exists a positive number $\alpha$ such that $\alpha \arg z \leq \arg F(z)$ for all $z \in \mathbb{H}_{+}, h(t)$ is operator monotone on $[0, \infty)$.

Proof. (1) Using the same argument of proof of Theorem 5 (1), it suffices to show that $0 \leq \arg h(z) \leq \pi$ for $z \in \mathbb{R}$ or $z \in \mathbb{H}_{+}$whose absolutely value is sufficiently large.

In the case $z \in(-\infty, 0)$, i.e., $|z|>0$ and $\arg z=\pi$, we have

$$
\begin{aligned}
& \quad \arg h(z) \\
& =\sum_{i=1}^{m} \arg \left(z-a_{i}\right)+\sum_{j=1}^{n} \arg \left(g_{j}(z)\left(z-b_{j}\right)\right) \\
& \quad-\sum_{i=1}^{n} \arg \left(f_{i}(z)-f_{i}\left(a_{i}\right)\right)-\sum_{j=1}^{n}\left(z g_{j}(z)-b_{j} g_{j}\left(b_{j}\right)\right) \\
& \leq m \pi+n \pi+\sum_{j=1}^{n} \arg g_{j}(z)-\sum_{i=1}^{m} \arg f_{i}(z)-n \pi \\
& = \\
& =\pi-\arg \frac{\prod_{i=1}^{n} f_{i}(z)}{z^{m-1} \prod_{j=1}^{n} g_{j}(z)} \leq \pi .
\end{aligned}
$$

So it holds $0 \leq \arg h(z) \leq \pi$.
In the case that $z \in \mathbb{H}_{+}$satisifies

$$
|z|>(m+n-1) \max \left\{a_{i}, b_{j} \mid 1 \leq i \leq m, 1 \leq j \leq n\right\} .
$$

Then it holds that

$$
\arg \left(z-a_{i}\right), \arg \left(z-b_{j}\right)<\frac{\pi+(m+n-1) \arg z}{m+n}
$$

by Lemma 4 . We may assume that there exists $k(1 \leq k \leq n)$ such that

$$
\arg \left(z g_{j}(z)\right) \leq \pi \quad(j \leq k), \quad \arg \left(z g_{j}(z)\right)>\pi \quad(j>k)
$$

Since

$$
\begin{aligned}
& \quad \arg h(z) \\
& \leq \frac{\pi+(m+n-1) \arg z}{m+n} \times m+\frac{\pi+(m+n-1) \arg z}{m+n} \times n+\sum_{j=1}^{n} \arg g_{j}(z) \\
& \quad-\sum_{i=1}^{m} \arg f_{i}(z)-\sum_{j=1}^{k} \arg z g_{j}(z)-(n-k) \pi \\
& =\pi+(m+n-k-1) \arg z+\sum_{j=k+1}^{n} \arg g_{j}(z)-\sum_{i=1}^{m} \arg f_{i}(z)-(n-k) \pi \\
& \leq \pi+(n-k)(\arg z-\pi)-\arg \frac{\prod_{i=1}^{m} f_{i}(z)}{z^{m-1} \prod_{j=1}^{m} g_{j}(z)} \\
& \leq \pi-\arg F(z) \leq \pi
\end{aligned}
$$

we have $0 \leq \arg h(z) \leq \pi$. So $h(t)$ is operator monotone on $[0, \infty)$.
(2) We choose a positive number $p$ as follows:

$$
\frac{m-1}{\alpha+m-1}<p<1
$$

We define functions $f_{i, p}, g_{j, p}$ as follows:

$$
f_{i, p}(z)=f_{i}\left(z^{p}\right), \quad g_{j, p}(z)=g_{j}\left(z^{p}\right) \quad\left(z \in \mathbb{H}_{+}\right) .
$$

Since $f_{i}, g_{j} \in \mathbb{P}_{+}[0, \infty), f_{i, p}, g_{j, p}$ are continuous on $\overline{\mathbb{H}_{+}}$and satisfy the condition

$$
f_{i, p}(t)-f_{i, p}\left(a_{i}\right) \neq 0, \quad t g_{j, p}(t)-b_{j} g_{j, p}\left(b_{j}\right) \neq 0, \quad t \in(-\infty, 0)
$$

by Lemma 3. We put

$$
F_{p}(t)=\frac{\prod_{i=1}^{m} f_{i, p}(t)}{t^{m-1} \prod_{j=1}^{n} g_{j, p}(t)}=F\left(t^{p}\right) t^{-(m-1)(1-p)}
$$

Then $F_{p}$ is holomorphic on $\mathbb{H}_{+}$and satisfies $F_{p}((0, \infty)) \subset(0, \infty)$. For any $z \in \mathbb{H}_{+}$, we have

$$
\arg F_{p}(z)=\arg F\left(z^{p}\right)-(m-1)(1-p) \arg z \leq \arg F\left(z^{p}\right) \leq \pi
$$

and

$$
\begin{aligned}
\arg F_{p}(z) & \geq \alpha \arg z^{p}-(m-1)(1-p) \arg z \\
& =(\alpha p-(m-1)(1-p)) \arg z \\
& =((\alpha+m-1) p-(m-1)) \arg z>0 .
\end{aligned}
$$

So we can see $F_{p} \in \mathbb{P}_{+}[0, \infty)$. By (1), we can show that

$$
h_{p}(t)=\prod_{i=1}^{m} \frac{\left(t-a_{i}\right)}{f_{i, p}(t)-f_{i, p}\left(a_{i}\right)} \prod_{j=1}^{n} \frac{g_{j, p}(t)\left(t-b_{j}\right)}{t g_{j, p}(t)-b_{j} g_{j, p}\left(b_{j}\right)}
$$

is operator monotone on $[0, \infty)$. When $p$ tends to $1, h_{p}(t)$ also tends to $h(t)$. Hence $h(t)$ is operator monotone on $[0, \infty)$.

Example 10. Let $0<p_{i} \leq 1(i=1,2, \ldots, m)$ and $0 \leq q_{j} \leq 1(j=$ $1,2, \ldots, n)$. We put

$$
f_{i}(t)=t^{p_{i}}, \quad g_{j}(t)=t^{q_{j}} \quad(t \geq 0)
$$

By the calculation

$$
F(t)=\frac{\prod_{i=1}^{m} f_{i}(t)}{t^{m-1} \prod_{j=1}^{n} g_{j}(t)}=t^{\sum_{i=1}^{m} p_{i}-\sum_{j=1}^{n} q_{j}-(m-1)},
$$

we have, for real numbers $a_{i}, b_{j} \geq 0$,

$$
h(t)=t^{\sum_{j=1}^{n} q_{j}} \frac{\left(t-a_{1}\right) \cdots\left(t-a_{m}\right)\left(t-b_{1}\right) \cdots\left(t-b_{n}\right)}{\left(t^{p_{1}}-a_{1}^{p_{1}}\right) \cdots\left(t^{p_{m}}-a_{m}^{p_{m}}\right)\left(t^{1+q_{1}}-b_{1}^{1+q_{1}}\right) \cdots\left(t^{1+q_{n}}-b_{n}^{1+q_{n}}\right)}
$$

is operator monotone on $[0, \infty)$ by Theorem 9 if it holds

$$
0 \leq \sum_{i=1}^{m} p_{i}-\sum_{j=1}^{n} q_{j}-(m-1) \leq 1
$$

i.e., $F(t)$ is operator monotone on $[0, \infty)$.

When $\sum_{i=1}^{m} p_{i}=\sum_{j=1}^{n} q_{j}+(m-1)$, we can see that

$$
h(t)=\frac{t^{\sum_{j=1}^{n} q_{j}}(t-1)^{m+n}}{\prod_{i=1}^{m}\left(t^{p_{i}}-1\right) \prod_{j=1}^{n}\left(t^{1+q_{j}}-1\right)}
$$

is operator monotone on $[0, \infty)$ and symmetric.
We can easily check that, if $h_{1}, h_{2} \in \mathbb{P}_{+}[0, \infty)$ are symmetric, then the functions

$$
\begin{aligned}
& f(t)=h_{1}(t)^{1 / p} h_{2}(t)^{1-1 / p} \quad(p>1) \\
& g(t)=\frac{t}{h_{1}(t)}
\end{aligned}
$$

are also operator monotone on $[0, \infty)$ and symmetric.
Combining these facts, for $r_{i}, s_{i}(i=1,2, \ldots, n)$ with

$$
\begin{gathered}
0<r_{1}, \ldots, r_{c} \leq 1, \quad 1 \leq r_{c+1}, \ldots, r_{n} \leq 2 \\
0<s_{1}, \ldots, s_{d} \leq 1, \\
1 \leq s_{d+1}, \ldots, s_{n} \leq 2 \\
\sum_{i=1}^{c} r_{i}=\sum_{i=c+1}^{n} r_{i}-1, \quad \sum_{i=1}^{d} s_{i}=\sum_{i=d+1}^{n} s_{j}-1
\end{gathered}
$$

we can see that the function

$$
h(t)=\sqrt{t^{\gamma} \prod_{i=1}^{n} \frac{r_{i}\left(t^{s_{i}}-1\right)}{s_{i}\left(t^{r_{i}}-1\right)}}
$$

is operator monotone on $[0, \infty)$ and symmetric with $h(1)=1$, where $\gamma=$ $1-c+d+\sum_{i=1}^{c} r_{i}-\sum_{i=1}^{d} s_{i}$.

By such a way, we can also construct from a symmetric operator monotone function to new one.

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# JI-DISTRIBUTIVE, DUALLY QUASI-DE MORGAN SEMI-HEYTING AND HEYTING ALGEBRAS 

HANAMANTAGOUDA P. SANKAPPANAVAR<br>Dedicated to Professor P.N. Shivakumar<br>A Great Humanitarian who changed the course of my life


#### Abstract

The variety DQD of dually quasi-De Morgan semiHeyting algebras and several of its subvarieties were investigated in the series [26] - [31]. In this paper we define and investigate a new subvariety JID of DQD, called "JI-distributive, dually quasi-De Morgan semi-Heyting algebras", defined by the identity: $x^{\prime} \vee(y \rightarrow z) \approx\left(x^{\prime} \vee y\right) \rightarrow\left(x^{\prime} \vee z\right)$, as well as the (closely related) variety DSt of dually Stone semi-Heyting algebras. Firstly, we prove that DSt and JID are discriminator varieties of level 1 and level 2 respectively. Secondly, we give a characterization of subdirectly irreducible algebras of the subvariety $\mathbf{J I D}_{\mathbf{1}}$ of $\mathbf{J I D}$ of level 1. As applications, we derive that the variety $\mathbf{J I D}_{\mathbf{1}}$ is the join of the variety DSt and the variety of De Morgan Boolean semiHeyting algebras, give a concrete description of the subdirectly irreducible algebras in the subvariety $\mathbf{J I D L}_{\mathbf{1}}$ of $\mathbf{J I D}_{\mathbf{1}}$ defined by the linear identity: $(x \rightarrow y) \vee(y \rightarrow x) \approx 1$, and deduce that the variety $\mathbf{J I D L}_{\mathbf{1}}$ is the join of the variety $\mathbf{D S t H C}$ generated by the dually Stone Heyting chains and the variety generated by the 4 -element De Morgan Boolean Heyting algebra. Furthermore, we present an explicit description of the lattice of subvarieties of $\mathbf{J I D L}_{\mathbf{1}}$ and equational bases for all subvarieties of $\mathbf{J I D L}_{\mathbf{1}}$. Finally, we prove that the amalgamation property holds for all subvarieties of DStHC.


## 1. Introduction

The De Morgan (strong) negation and the pseudocomplement are two of the fairly well known negations that generalize the classical negation. A common generalization of these two negations led to a

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new variety of algebras, called "semi-De Morgan algebras", which was investigated in [24]. Several subvarieties of this variety, including a subvariety called "(upper) quasi-De Morgan algebras" were also studied in [24].

In a different vein, semi-Heyting algebras were introduced in [25] as an abstraction of Heyting algebras. Using the dual version of quasi-De Morgan negation, an expansion of semi-Heyting algebras, called "dually quasi-De Morgan semi-Heyting algebras (DQD, for short)" was defined and investigated in [26], as a common generalization of De Morgan (or symmetric) Heyting algebras [23] (see also [19]) and dually pseudocomplemented Heyting algebras [22]. It may also be mentioned here that [8] has proposed recently a propositional logic, called "Dually quasi-De Morgan semi-Heyting logic", which has dually quasi-De Morgan semi-Heyting algebras as an equivalent algebraic semantics.

Several new subvarieties of DQD were studied in [26]- [31], including the variety DStHC generated by the dually Stone Heyting chains (i.e., the expansion of the Gödel variety by the dual Stone operation), the variety DMB of De Morgan Boolean semi-Heyting algebras and the variety DMBH generated by the 4 -element De Morgan Boolean Heyting algebra. These investigations led us naturally to the problem of equational axiomatization for the join of the variety DStHC and the variety DMBH. Our investigations into this problem led us to the results of the present paper that include a solution to the just mentioned problem.

In this paper we define and investigate a new subvariety of DQD, called "JI-distributive, dually quasi-De Morgan semi-Heyting algebras (JID, for short)", defined by the identity: $x^{\prime} \vee(y \rightarrow z) \approx\left(x^{\prime} \vee y\right) \rightarrow$ $\left(x^{\prime} \vee z\right)$, as well as the (closely related) variety DSt of dually Stone semiHeyting algebras. We first prove that DSt and JID are discriminator varieties of level 1 and level 2 respectively (see Section 2 for definitions). Secondly, we prove that the lattice of subvarieties of DStHC is an $\omega+1$-chain. Thirdly, we give a characterization of subdirectly irreducible algebras of the subvariety $\mathbf{J I D}_{1}$ of level 1 . As a first application of it, we derive that the variety $\mathbf{J I D}_{\mathbf{1}}$ is the join of the variety $\mathbf{D S t}$ and the variety DMB. As a second application, we give a concrete description of the subdirectly irreducible algebras in the subvariety $\mathrm{JIDL}_{1}$ of $\mathbf{J I D}_{1}$ defined by the linear identity: $(x \rightarrow y) \vee(y \rightarrow x) \approx 1$, and deduce that the variety $\mathbf{J I D L}_{1}$ is the join of the variety $\mathbf{D S t H C}$ generated by the dually Stone Heyting chains and the variety DMBH. Other applications include a description of the lattice of subvarieties of $\mathbf{J I D L} \mathbf{1}_{1}$, equational bases of all subvarieties of $\mathbf{J I D L}_{\mathbf{1}}$, and the fact that the amalgamation property holds in all subvarieties of DStHC.

More explicitly, the paper is organized as follows: In Section 2 we recall definitions, notations and results from [26], [27] and [28] and also prove some new results needed in the rest of the paper. In Section 3, we define the variety JID of JI-distributive, dually quasi-De Morgan semi-Heyting algebras and give some arithmetical properties of JID. In particular, we show that JID satisfies the $\vee$-De Morgan law and the level 2 identity: $\left(x \wedge x^{\prime *}\right)^{\prime *} \approx\left(x \wedge x^{\prime *}\right)^{1 * * *}$. These two propertes allow us to apply [26, Corollary 8.2(a)] to deduce that JID is a discriminator variety. These properties also play a crucial role in the rest of the paper. Section 4 will prove that the variety $\mathbf{D S t}$ is a discriminator variety of level 1. It will also present some properties of DSt, which, besides being of interest in their own right, will also be useful in the later sections. It is also proved that the lattice of subvarieties of DStHC is an $\omega+1$-chain. In Section 5 , we give a characterization of subdirectly irreducible ( $=$ simple) algebras in the variety $\mathbf{J I D}_{\mathbf{1}}$ of level 1 and deduce that $\mathbf{J I D}_{\mathbf{1}}$ is the join of $\mathbf{D S t}$ and the variety $\mathbf{D M B}$ of De Morgan Boolean semi-Heyting algebras. Several applications of this characterization are given in Section 6 and Section 7. We investigate, in Section 6 , the variety $\mathrm{JIDL}_{1}$ of JI-distributive, dually quasi-De Morgan, linear semi-Heyting algebras of level 1. An explicit description of subdirectly irreducible algebras in $\mathbf{J I D L}_{\mathbf{1}}$ is given, and from this description it is deduced that $\mathbf{J I D L}_{1}=\mathbf{D S t H C} \vee \mathbf{D M B H}$, which solves the aforementioned problem of axiomatizing the join of DStHC and DMBH. In Section 7, some applications of the just-mentioned result are given. It is shown that the lattice of subvarieties of $\mathbf{J I D L}_{1}$ is isomorphic to $\mathbf{1} \oplus[(\omega+\mathbf{1}) \times \mathbf{2}]$, where $\mathbf{1}$ and $\mathbf{2}$ are the 1 -element and the 2 -element lattices, respectively. Also, (small) equational bases for all subvarieties of $\mathbf{J I D L}_{\mathbf{1}}$ are given. Finally, it is shown that all subvarieties of $\mathbf{D S t H C}$ have the amalgamation property.

## 2. Preliminaries

In this section we recall some notions and known results needed to make this paper as self-contained as possible. However, for other information used but not mentioned here, we refer the reader to [5], [7] and [20].

An algebra $\mathbf{L}=\langle L, \vee, \wedge, \rightarrow, 0,1\rangle$ is a semi-Heyting algebra ([25]) if $\langle L, \vee, \wedge, 0,1\rangle$ is a bounded lattice and $\mathbf{L}$ satisfies:
(SH1) $x \wedge(x \rightarrow y) \approx x \wedge y$,
(SH2) $x \wedge(y \rightarrow z) \approx x \wedge[(x \wedge y) \rightarrow(x \wedge z)]$,
(SH3) $x \rightarrow x \approx 1$.
Semi-Heyting algebras are distributive and pseudocomplemented, with $a^{*}:=a \rightarrow 0$ as the pseudocomplement of an element $a$.
Let $\mathbf{L}$ be a semi-Heyting algebra. $\mathbf{L}$ is a Heyting algebra if $\mathbf{L}$ satisfies:
(H) $(x \wedge y) \rightarrow y \approx 1$.
$\mathbf{L}$ is a Boolean semi-Heyting algebra if $\mathbf{L}$ satisfies:
(Bo) $x \vee x^{*} \approx 1$.
$\mathbf{L}$ is a Boolean Heyting algebra if $\mathbf{L}$ is a Heyting algebra and satisfies (Bo).

The following definition, taken from [26], is central to this paper.
DEFINITION 2.1. An algebra $\mathbf{L}=\left\langle L, \vee, \wedge, \rightarrow,{ }^{\prime}, 0,1\right\rangle$ is a semiHeyting algebra with a dual quasi-De Morgan operation or dually quasiDe Morgan semi-Heyting algebra (DQD-algebra, for short) if $\langle L, \vee, \wedge, \rightarrow, 0,1\rangle$ is a semi-Heyting algebra, and $\mathbf{L}$ satisfies:
(a) $0^{\prime} \approx 1$ and $1^{\prime} \approx 0$,
(b) $(x \wedge y)^{\prime} \approx x^{\prime} \vee y^{\prime}$,
(c) $(x \vee y)^{\prime \prime} \approx x^{\prime \prime} \vee y^{\prime \prime}$,
(d) $x^{\prime \prime} \leq x$.

Let $\mathbf{L}$ be a DQD-algebra. $\mathbf{L}$ is a dually pseudocomplemented semiHeyting algebra (DPC-algebra) (see [24]) if $\mathbf{L}$ satisfies:
(e) $x \vee x^{\prime} \approx 1$.
$\mathbf{L}$ is a dually Stone semi-Heyting algebra (DSt-algebra) if $\mathbf{L}$ satisfies the dual Stone identity:
$(\mathrm{DSt}) x^{\prime} \wedge x^{\prime \prime} \approx 0$.
It should be noted that if (DSt) holds in a DQD-algebra $\mathbf{L}$, then (e) holds in $\mathbf{L}$ as well, and hence ' is indeed the dual pseudocomplement satisfying the dual Stone identity, and so $\mathbf{L}$ has, indeed, a dual Stone algebra as a reduct. $\mathbf{L}$ is a De Morgan semi-Heyting algebra (DMalgebra) if $\mathbf{L}$ satisfies:
$(\mathrm{DM}) x^{\prime \prime} \approx x$.
The varieties of DQD-algebras, DPC-algebras, DSt-algebras, DMalgebras are denoted, respectively, by DQD, DPC, DSt, and DM. If the underlying semi-Heyting algebra of a DQD-algebra is a Heyting algebra, then we add "H" at the end of the names of the varieties that will be considered in the sequel. Thus, for example, DStH denotes the variety of dually Stone Heyting algebras.

The following lemmas are basic to this paper. The proof of the first lemma is straightforward and is left to the reader.

LEMMA 2.2. Let $\mathbf{L} \in \mathrm{DQD}$ and let $x, y, z \in L$. Then
(i) $1^{\prime *}=1$, and $1 \rightarrow x=x$,
(ii) $x \leq y$ implies $x^{\prime} \geq y^{\prime}$,
(iii) $(x \wedge y)^{\prime *}=x^{*} \wedge y^{\prime *}$,
(iv) $x^{\prime \prime \prime}=x^{\prime}$,
(v) $(x \vee y)^{\prime}=\left(x^{\prime \prime} \vee y\right)^{\prime}$,
(vi) $x \wedge[y \vee(x \rightarrow z)]=x \wedge(y \vee z)$,
(vii) $x \wedge(x \rightarrow y)^{\prime \prime} \leq y$.

LEMMA 2.3. Let $L \in \mathrm{DQD}$ and $x, y \in L$. Then
(1) $(x \vee y)^{\prime} \leq x^{\prime} \rightarrow(x \vee y)^{\prime}$,
(2) $\left[x \vee(y \vee z)^{\prime}\right]^{\prime}=\left(x \vee y^{\prime}\right)^{\prime} \vee\left(x \vee z^{\prime}\right)^{\prime}$,
(3) $x \wedge[(x \rightarrow y) \vee z]=x \wedge(y \vee z)$,
(4) $y \wedge[x \rightarrow(y \wedge z)]=y \wedge(x \rightarrow z)$,
(5) $x \rightarrow(y \wedge z) \geq y \wedge(x \rightarrow z)$,
(6) $x \leq y \rightarrow(x \wedge y)$,
(7) $(x \vee y)^{\prime}=x^{\prime} \wedge\left[(x \vee y)^{\prime} \vee\left\{x^{\prime} \rightarrow(x \vee y)^{\prime}\right\}^{\prime \prime}\right]$,
(8) $x \leq(x \rightarrow y) \rightarrow y$.

Proof.
(1) is straightforward to verify since $(x \vee y)^{\prime} \leq x^{\prime}$.
(2): $\left[\left(x \vee(y \vee z)^{\prime}\right]^{\prime}=\left[x^{\prime \prime} \vee(y \vee z)^{\prime}\right]^{\prime} \quad\right.$ by Lemma $2.2(\mathrm{v})$

$$
=\left[x^{\prime} \wedge(y \vee z)\right]^{\prime \prime}
$$

$$
=\left[\left(x^{\prime} \wedge y\right) \vee\left(x^{\prime} \wedge z\right)\right]^{\prime \prime}
$$

$$
=\left(x^{\prime} \wedge y\right)^{\prime \prime} \vee\left(x^{\prime} \wedge z\right)^{\prime \prime}
$$

$$
=\left(x^{\prime \prime} \vee y^{\prime}\right)^{\prime} \vee\left(x^{\prime \prime} \vee z^{\prime}\right)^{\prime}
$$

$$
=\left(x \vee y^{\prime}\right)^{\prime} \vee\left(x \vee z^{\prime}\right)^{\prime} \quad \text { by Lemma } 2.2(\mathrm{v})
$$

(3) and (4) are easy to verify.
(5): $[x \rightarrow(y \wedge z)] \wedge y \wedge(x \rightarrow z)=y \wedge[x \rightarrow(y \wedge z)] \wedge(x \rightarrow z)$ $=y \wedge(x \rightarrow z)$ by (4).
(6): $x=x \wedge(y \rightarrow y) \leq y \rightarrow(x \wedge y)$ by (5).

$$
\begin{aligned}
(7):(x \vee y)^{\prime} & =x^{\prime} \wedge(x \vee y)^{\prime} \\
& =x^{\prime} \wedge\left[x^{\prime} \rightarrow(x \vee y)^{\prime}\right] \\
& =x^{\prime} \wedge\left[\left\{x^{\prime} \rightarrow(x \vee y)^{\prime}\right\} \vee\left\{x^{\prime} \rightarrow(x \vee y)^{\prime}\right\}^{\prime \prime}\right] \\
& =\left[x^{\prime} \wedge(x \vee y)^{\prime}\right] \vee\left[x^{\prime} \wedge\left\{x^{\prime} \rightarrow(x \vee y)^{\prime}\right\}^{\prime \prime}\right] \\
& =x^{\prime} \wedge\left[(x \vee y)^{\prime} \vee\left\{x^{\prime} \rightarrow(x \vee y)^{\prime}\right\}^{\prime \prime}\right],
\end{aligned}
$$

(8): $x \wedge[(x \rightarrow y) \rightarrow y]=x \wedge[\{x \wedge(x \rightarrow y)\} \rightarrow(x \wedge y)]=x \wedge[(x \wedge y) \rightarrow$ $(x \wedge y)]=x \wedge 1=x$, completing the proof.

The following three 4-element algebras, called $\mathbf{D}_{1}, \mathbf{D}_{\mathbf{2}}$, and $\mathbf{D}_{3}$ (following the notation of [26]), in DQD, play an important role in the sequel. All three of them have the Boolean lattice reduct with the universe $\{0, a, b, 1\}$, where $b$ is the Boolean complement of $a$, and the operation ' is defined as follows: $a^{\prime}=a, b^{\prime}=b, 0^{\prime}=1,1^{\prime}=0$, while the operation $\rightarrow$ is defined in Figure 1.

| $\rightarrow$ | 0 | 1 | $a$ | $b$ |
| :---: | :---: | :---: | :---: | :---: |
| 0 | 1 | 0 | $b$ | $a$ |
| $\mathrm{D}_{1}: 1$ | 0 | 1 | $a$ | $b$ |
| $a$ | $b$ | $a$ | 1 | 0 |
| $b$ | $a$ | $b$ | 0 | 1 |


| $\rightarrow$ | 0 | 1 | $a$ | $b$ |
| ---: | ---: | ---: | ---: | ---: |
| 0 | 1 | 1 | 1 | 1 |
| $: 1$ | 0 | 1 | $a$ | $b$ |
| a | $b$ | 1 | 1 | $b$ |
| $b$ | $a$ | 1 | $a$ | 1 |


$\mathbf{D}_{\mathbf{3}}$|  |
| :---: |
|  | |  | $\rightarrow$ | 0 | 1 | $a$ |
| :---: | :---: | :---: | :---: | :---: |
| 0 | $b$ |  |  |  |
|  | 1 | 1 | $a$ | 1 |
| $a$ |  |  |  |  |
|  | 0 | 1 | $a$ | $b$ |
| $a$ | $b$ | $a$ | 1 | 0 |
| $b$ | $a$ | 1 | $a$ | 1 |

Figure 1

Let DQB and DMB denote respectively the subvarieties of DQD and DM defined by (Bo). Also, by an earlier convention, DQBH and DMBH denote, respectively, the subvarieties of DQB and DMB defined by $(\mathrm{H}) . \mathbf{V}(\mathbf{K})$ denotes the variety generated by the class K of algebras in DQD. The following proposition is proved in ([26]) and is needed later in this paper.

## PROPOSITION 2.4.

(a) $\mathrm{DQB}=\mathbf{D M B}=\mathbf{V}\left(\left\{\mathbf{D}_{\mathbf{1}}, \mathbf{D}_{\mathbf{2}}, \mathbf{D}_{\mathbf{3}}\right\}\right)$,
(b) $\mathrm{DQBH}=\mathrm{DMBH}=\mathrm{V}\left(\mathrm{D}_{2}\right)$.

The following definition is from [26].
DEFINITION 2.5. Let $\mathbf{L} \in \mathbf{D Q D}$ and $x \in \mathbf{L}$. For $n \in \omega$, we define $t_{n}(x)$ recursively as follows:

$$
\begin{gathered}
x^{0(1 *)}:=x \\
x^{(n+1)(/ *)}:=\left(x^{n(/ *))^{\prime *}}, \text { for } n \geq 0\right. \\
t_{0}(x):=x \\
t_{n+1}(x):=t_{n}(x) \wedge x^{(n+1)(/ *)}, \text { for } n \geq 0
\end{gathered}
$$

Let $n \in \omega$. The subvariety $\mathbf{D Q D}_{\mathbf{n}}$ of level $n$ of $\mathbf{D Q D}$ is defined by the identity:
(lev $n$ )

$$
t_{n}(x) \approx t_{(n+1)}(x)
$$

For a subvariety $\mathbf{V}$ of $\mathbf{D Q D}$, we let $\mathbf{V}_{\mathbf{n}}:=\mathbf{V} \cap \mathbf{D Q D}_{\mathbf{n}}$.
Recall from [26] (or [27]) that BDQDSH is the subvariety of DQD (= DQDSH) defined by the identity:
(BL) $\left(x \vee y^{*}\right)^{\prime} \approx x^{\prime} \wedge y^{* \prime}$.
We will abbreviate BDQDSH by BDQD.
The following "simplicity condition", (SC), is crucial in the rest of the paper.
(SC) For every $x \in L$, if $x \neq 1$, then $x \wedge x^{\prime *}=0$.
The following theorem, which was proved in [27, Corollary 4.1] (which is, in turn, a consequence of Corollaries 7.6 and 7.7 of [26]), will play a fundamental role in this paper.

THEOREM 2.6. [27, Corollary 4.1] Let $\mathbf{L} \in \mathbf{B D Q D}_{1}$ with $|L| \geq 2$. Then the following are equivalent:
(1) $\mathbf{L}$ is simple,
(2) $\mathbf{L}$ is subdirectly irreducible,
(3) L satisfies (SC).

## 3. JI-distributive, dually quasi-De Morgan semi-Heyting ALGEBRAS

The identity, $x \vee(y \rightarrow z) \approx(x \vee y) \rightarrow(x \vee z)$, was shown in [28, Corollary 3.55] to be an equational base for the variety generated by $\mathbf{D}_{\mathbf{2}}$, relative to $\mathbf{D Q D}$. Let us refer to this identity as "strong

JI-distributive identity". We now introduce a slightly weaker identity, called "JI-distributive identity" (by restricting the first variable to "primed" elements). The subvariety JID of DQD defined by this identity and some of its subvarieties are the subject of our investigation in the rest of this paper.
DEFINITION 3.1. The subvariety JID of DQD is defined by: (JID) $x^{\prime} \vee(y \rightarrow z) \approx\left(x^{\prime} \vee y\right) \rightarrow\left(x^{\prime} \vee z\right)$ (restricted Distribution of Join over Implication).
Members of the variety JID are called "JI-distributive, dually quasiDe Morgan semi-Heyting algebras" and will be referred to as JIDalgebras. Examples of JID-algebras come from a surprising source to which we shall now turn. But, first we need some notation.

A DQD-algebra is a DQD-chain if its lattice reduct is a chain. Let DQDC [DPCC] denote the variety generated by the DQD-chains [DPC-chains]. The following lemma provides an important class of examples of JID, which is partly the motivation for our interest in JID.

## LEMMA 3.2. DPCC $\subseteq$ JID.

Proof. It suffice to show that DPCC $\models$ (JID). Let A be a DPC-chain and let $a \in A \backslash\{1\}$. Since A is a chain, we have $a^{\prime} \leq a$ or $a \leq a^{\prime}$, from which we get that $a \vee a^{\prime} \leq a$ or $a \vee a^{\prime} \leq a^{\prime}$. Since $\mathbf{A}$ is dually pseudocomplemented, we have $a \vee a^{\prime}=1$, implying $a^{\prime}=1$, as $a \neq 1$. Now, it is routine to verify (JID) holds in A.

For $\mathbf{L}$ a DPC-chain, it was observed in the proof of the preceding lemma that the dual pseudocomplement ' satisfies: $a^{\prime}=1$, if $a \neq 1$, and hence $\mathbf{L} \models(\mathrm{DSt})$. Thus, we have the following corollary, where DStC denotes the variety generated by the dually Stone semi-Heyting chains.

COROLLARY 3.3. $\mathrm{DPCC}=\mathrm{DStC}$.
From now on, we use DPCC and DStC interchangeably. We note that $\mathbf{D}_{\mathbf{1}}, \mathbf{D}_{\mathbf{2}}$, and $\mathbf{D}_{\mathbf{3}}$ are also examples of JID-algebras.

In the rest of this section we present several useful arithmetical properties of JID. Following our convention made earlier, JIDH denotes the subvariety of JID defined by the identity (H).

Throughout this section, we assume that $\mathrm{L} \in \mathrm{JID}$.
LEMMA 3.4. Let $x, y, z \in \mathbf{L}$. Then
(1) $x^{\prime} \rightarrow\left(x^{\prime} \vee y\right)=x^{\prime} \vee\left(x^{\prime} \rightarrow y\right)$,
(2) $x^{\prime} \rightarrow\left(x^{\prime} \vee y\right)=x^{\prime} \vee(0 \rightarrow y)$,
(3) $x^{\prime} \vee\left(x^{\prime} \rightarrow y\right)=x^{\prime} \vee(0 \rightarrow y)$; in particular, $x^{\prime} \vee x^{\prime *}=1$,
(4) $\left(x^{\prime} \vee y\right) \rightarrow x^{\prime}=x^{\prime} \vee y^{*}$,
(5) $\left(x^{\prime} \vee y\right) \rightarrow x^{\prime}=x^{\prime} \vee\left(y \rightarrow x^{\prime}\right)$,
(6) $x^{\prime} \vee\left(y \rightarrow x^{\prime}\right)=x^{\prime} \vee y^{*}$,
(7) $x^{\prime} \rightarrow(x \vee y)^{\prime}=x^{*} \vee(x \vee y)^{\prime}$.

Proof. Observe that $x^{\prime} \rightarrow\left(x^{\prime} \vee y\right)=\left(x^{\prime} \vee x^{\prime}\right) \rightarrow\left(x^{\prime} \vee y\right)=x^{\prime} \vee\left(x^{\prime} \rightarrow y\right)$ by (JID), which proves (1). To prove (2), again using (JID), we get $x^{\prime} \vee(0 \rightarrow y)=\left(x^{\prime} \vee 0\right) \rightarrow\left(x^{\prime} \vee y\right)=x^{\prime} \rightarrow\left(x^{\prime} \vee y\right)$. (3) is immediate from (1) and (2). For (4), $\left(x^{\prime} \vee y\right) \rightarrow x^{\prime}=\left(x^{\prime} \vee y\right) \rightarrow\left(x^{\prime} \vee 0\right)=x^{\prime} \vee(y \rightarrow 0)=$ $x^{\prime} \vee y^{*}$, in view of (JID). Next, $\left(x^{\prime} \vee y\right) \rightarrow x^{\prime}=\left(x^{\prime} \vee y\right) \rightarrow\left(x^{\prime} \vee x^{\prime}\right)=$ $x^{\prime} \vee\left(y \rightarrow x^{\prime}\right)$, proving (5), and (6) is immediate from (4) and (5). For (7), we have

$$
\begin{aligned}
x^{\prime} \rightarrow(x \vee y)^{\prime} & =(x \vee y)^{\prime} \vee\left[x^{\prime} \rightarrow(x \vee y)^{\prime}\right] \text { by Lemma } 2.3 \\
& =(x \vee y)^{\prime} \vee x^{\prime *} \text { by }(6) .
\end{aligned}
$$

We now prove an important property of the variety JID, namely the V-De Morgan law. We denote by Dms the subvariety of DQD (called "dually ms semi-Heyting algebras") defined by

$$
(x \vee y)^{\prime} \approx x^{\prime} \wedge y^{\prime} \quad(\vee-\text { De Morgan Law })
$$

## THEOREM 3.5. JID $\subseteq$ Dms.

Proof. Let $x, y \in \mathbf{L}$. As $x^{\prime} \wedge x^{\prime * \prime \prime} \leq x^{\prime} \wedge x^{\prime *}=0$, we get $x^{\prime} \wedge y^{\prime}=$ $\left(x^{\prime} \wedge x^{\prime * \prime \prime}\right) \vee\left(x^{\prime} \wedge y^{\prime}\right)$. Hence,

$$
\begin{aligned}
x^{\prime} \wedge y^{\prime} & =x^{\prime} \wedge\left(x^{\prime * \prime \prime} \vee y^{\prime}\right) & & \\
& =x^{\prime} \wedge\left[(x \vee y)^{\prime} \vee x^{\prime * \prime \prime} \vee y^{\prime}\right] & & \text { since }(x \vee y)^{\prime} \leq y^{\prime} \\
& =x^{\prime} \wedge\left[(x \vee y)^{\prime} \vee x^{\prime * \prime \prime} \vee y^{\prime \prime \prime}\right] & & \\
& =x^{\prime} \wedge\left[(x \vee y)^{\prime} \vee\left(x^{\prime *} \vee y^{\prime}\right)^{\prime \prime}\right] & & \text { by Lemma 3.4 (3.4) } \\
& =x^{\prime} \wedge\left[(x \vee y)^{\prime} \vee\left\{\left(x^{\prime *} \vee x^{\prime}\right)^{\prime} \vee\left(x^{\prime *} \vee y^{\prime}\right)^{\prime}\right\}^{\prime}\right] & & \text { by Lemma 2.3 (2) } \\
& =x^{\prime} \wedge\left[(x \vee y)^{\prime} \vee\left\{x^{\prime *} \vee(x \vee y)^{\prime}\right\}^{\prime \prime}\right] & & \text { by Lemma 3.4 (7) } \\
& =x^{\prime} \wedge\left[(x \vee y)^{\prime} \vee\left\{x^{\prime} \rightarrow(x \vee y)^{\prime}\right\}^{\prime \prime \prime}\right] & & \text { by Lemma 2.3 (7). }
\end{aligned}
$$

Hence, JID $\subseteq$ Dms.
The following lemma is useful in this and later sections.
LEMMA 3.6. Let $x, y, z \in \mathbf{L}$. Then
(1) $x^{\prime * \prime \prime}=x^{\prime *}$,
(2) $x^{\prime \prime *}=x^{\prime * \prime}$,
(3) $x \rightarrow\left(x \wedge y^{\prime}\right)=x^{*} \vee y^{\prime}$,
(4) $\left(x \wedge y^{\prime *}\right)^{*}=y^{\prime} \vee x^{*}$,
(5) $\left(x^{\prime} \vee y^{\prime \prime *}\right)^{* \prime}=\left(x^{\prime \prime} \wedge y^{\prime *}\right)^{*}$.

Proof. (1): From Lemma 3.4 (3) we have $x^{\prime} \vee x^{*}=1$, which yields $x^{\prime \prime \prime} \vee x^{\prime * \prime \prime}=1$, implying $x^{\prime} \vee x^{\prime * \prime \prime}=1$, leading to $x^{\prime *} \leq x^{\prime * \prime \prime}$; thus, $x^{\prime *}=x^{\prime * \prime \prime}$.
(2): From $x^{\prime} \vee x^{\prime *}=1$ and Theorem 3.5 we get $x^{\prime \prime} \wedge x^{\prime * \prime}=0$, implying $x^{\prime * \prime} \leq x^{\prime \prime *}$. To prove the reverse inequality, from $x^{\prime} \wedge x^{\prime *}=0$, we get $x^{\prime \prime} \vee x^{\prime * \prime}=1$, from which it follows that $x^{\prime * *} \leq x^{\prime * \prime}$.

$$
\begin{aligned}
& \text { (3): } x^{*} \vee y^{\prime}=\left(y^{\prime} \vee x\right) \rightarrow y^{\prime} \\
& \text { by (JID) } \\
& =\left(y^{\prime} \vee x\right) \rightarrow\left[y^{\prime} \vee\left(x \wedge y^{\prime}\right)\right] \\
& =y^{\prime} \vee\left[x \rightarrow\left(x \wedge y^{\prime}\right)\right] \quad \text { by }(\mathrm{JID}) \\
& =x \rightarrow\left(x \wedge y^{\prime}\right) \quad \text { by Lemma } 2.3 \text { (6). } \\
& \text { (4): }\left(x \wedge y^{\prime *}\right)^{*}=\left(x \wedge y^{\prime \prime \prime *}\right)^{*} \\
& =\left(x \wedge y^{\prime * \prime \prime}\right)^{*} \quad \text { by (2) (twice) } \\
& =\left(x \wedge y^{\prime * \prime \prime}\right) \rightarrow\left(y^{\prime} \wedge x \wedge y^{\prime * \prime \prime}\right) \quad \text { as } y^{\prime} \wedge y^{\prime * \prime}=0 \\
& =y^{\prime} \vee\left(x \wedge y^{\prime * \prime \prime}\right)^{*} \quad \text { by (3) } \\
& =y^{\prime} \vee\left(x \wedge y^{\prime \prime \prime}\right)^{*} \quad \text { by (2) (twice) } \\
& =y^{\prime} \vee\left(x \wedge y^{\prime *}\right)^{*} \\
& =y^{\prime} \vee\left[\left(x \wedge y^{\prime *}\right) \rightarrow 0\right] \\
& =\left[y^{\prime} \vee\left(x \wedge y^{\prime *}\right)\right] \rightarrow y^{\prime} \quad \text { by (JID) } \\
& =\left[\left(y^{\prime} \vee x\right) \wedge\left(y^{\prime} \vee y^{\prime *}\right)\right] \rightarrow y^{\prime} \\
& =\left[\left(y^{\prime} \vee x\right) \wedge 1\right] \rightarrow y^{\prime} \quad \text { by Lemma } 3.4 \text { (3) } \\
& =\left(y^{\prime} \vee x\right) \rightarrow y^{\prime} \\
& =y^{\prime} \vee(x \rightarrow 0) \quad \text { by (JID) } \\
& =y^{\prime} \vee x^{*} \text {. }
\end{aligned}
$$

$$
\begin{array}{rlr}
(5):\left(x^{\prime} \vee y^{\prime \prime *}\right)^{* \prime} & =\left(x^{\prime} \vee y^{\prime * \prime}\right)^{* \prime} & \text { by }(2) \\
& =\left(x \wedge y^{\prime *}\right)^{\prime * \prime} & \\
& =\left(x \wedge y^{\prime *}\right)^{\prime \prime *} & \text { by }(2) \\
& =\left(x^{\prime} \vee y^{\prime * \prime}\right)^{\prime *} & \\
& =\left(x^{\prime \prime} \wedge y^{\prime * \prime \prime}\right)^{*} & \text { by Theorem } 3.5 \\
& =\left(x^{\prime \prime} \wedge y^{\prime \prime \prime *}\right)^{*} & \text { by }(2)(\text { twice }) \\
& =\left(x^{\prime \prime} \wedge y^{\prime *}\right)^{*} . &
\end{array}
$$

This completes the proof.

### 3.1. An Alternate Definition of "level $n$ ", for $n \geq 1$.

The following lemmas enable us to give an alternate definition of "Level $n$ ", for $n \geq 1$.

LEMMA 3.7. Let $x \in \mathbf{L}$. Then $x^{* * *}=x^{\prime}$.
Proof. Since $x^{\prime} \vee x^{* *}=1$ by Lemma 3.4, and $x^{\prime} \wedge x^{* *}=0$, we get $x^{\prime * *}=x^{\prime}$.

LEMMA 3.8. Let $x \in \mathbf{L}$. Then $x \wedge x^{* *} \wedge x^{\prime * * *}=\left(x \wedge x^{\prime *}\right)^{\prime *}$.

$$
\text { Proof. } \begin{array}{rlr}
x \wedge x^{\prime *} \wedge x^{\prime * * *} & =x \wedge x^{\prime *} \wedge x^{\prime \prime * *} & \text { by Lemma } 3.6(2) \\
& =x \wedge x^{\prime *} \wedge x^{\prime \prime} & \\
& =x^{\prime *} \wedge x^{\prime \prime} & \\
& =x^{\prime *} \wedge x^{\prime \prime * *} & \\
& =x^{\prime *} \wedge x^{\prime * *} & \text { by Lemma } 3.7 \\
& =\left(x \wedge x^{\prime *}\right)^{\prime *} . &
\end{array}
$$

Since $\mathbf{J I D}_{\mathbf{n}}=\mathbf{J I D} \cap \mathbf{D Q D}_{\mathbf{n}}$, the above lemma allows us to make the following alternate (but equivalent) definition for $\mathbf{J I D}_{\mathbf{n}}$, for $n \in \omega$ such that $n \geq 1$.

DEFINITION 3.9. Let $n$ be an integer $\geq 1$. The variety $\mathbf{J I D}_{\mathbf{n}}$ is the subvariety of JID defined by
$\left(\right.$ Lev n) $\quad\left(x \wedge x^{\prime *}\right)^{(n-1)(1 *)} \approx\left(x \wedge x^{*}\right)^{n(1 *)}$.
Thus, in particular, $\mathbf{J I D}_{\mathbf{1}}$ and $\mathbf{J I D}_{\mathbf{2}}$ are, respectively, defined, relative to JID, by
(Lev 1) $x \wedge x^{* *} \approx\left(x \wedge x^{*}\right)^{* *}$,
(Lev 2) $\quad\left(x \wedge x^{* *}\right)^{* *} \approx\left(x \wedge x^{* *}\right)^{* * *}$.

In the rest of the paper we will use these definitions for the levels of $\mathrm{JID}_{1}$ and $\mathbf{J I D}_{2}$.

### 3.2. The Level of JID.

Next, we wish to prove that JID is at Level 2.
THEOREM 3.10. We have
(1) $\mathrm{JID}_{\mathbf{1}} \subset \mathrm{JID}$,
(2) $\mathbf{J I D}=\mathbf{J I D}_{\mathbf{2}}$.

Proof. First, we prove (2). That is, we need to prove that the "level 2" identity holds in JID. Let $x \in \mathbf{L}$.

$$
\begin{aligned}
\left(x \wedge x^{\prime *}\right)^{\prime * / *} & =\left(x^{\prime} \vee x^{\prime * \prime}\right)^{* / *} & & \\
& =\left(x^{\prime} \vee x^{\prime \prime *}\right)^{* / *} & & \text { by Lemma } 3.6(2) \\
& =\left(x^{\prime \prime} \wedge x^{\prime *}\right)^{* *} & & \text { by Lemma } 3.6(5) \\
& =\left(x^{\prime} \vee x^{\prime \prime *}\right)^{*} & & \text { by Lemma 3.6 (4) } \\
& =\left(x^{\prime} \vee x^{\prime * \prime}\right)^{*} & & \text { by Lemma 3.6 (2) } \\
& =\left(x \wedge x^{\prime *}\right)^{\prime *} . & &
\end{aligned}
$$

Hence (2) is proved. For (1), we consider the following algebra SIX with its lattice reduct, $\rightarrow$ and ' as given in Figure 2. We note that SIX $\in \mathbf{J I D}$; but it is not of level 1 (at a).


| $\rightarrow:$ | 0 | 1 | $a$ | $b$ | $c$ | $d$ |  |  |  |  |  |  |
| ---: | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| 0 | 1 | 1 | 1 | 1 | 1 | 1 |  |  |  |  |  |  |
| 1 | 0 | 1 | $a$ | $b$ | $c$ | $d$ |  |  |  |  |  |  |
| $a$ | 0 | 1 | 1 | $b$ | $c$ | $c$ |  |  |  |  |  |  |
| $b$ | $c$ | 1 | 1 | 1 | $c$ | $c$ |  |  |  |  |  |  |
| $c$ |  |  | $:$ |  | 0 | 1 | $a$ | $b$ | $c$ | $d$ |  |  |
|  |  |  | 1 | 1 | 0 | $b$ | $b$ | $c$ | 1 |  |  |  |
| $d$ | $b$ | 1 | 1 | $b$ | 1 | $a$ |  |  |  |  |  |  |
|  |  |  |  |  |  |  |  |  |  |  |  |  |

Figure 2

The following corollary is immediate from the above theorem and [26, Corollary 8.2(a)].

COROLLARY 3.11. JID is a discriminator variety of level 2.

## 4. Dually Stone Semi-Heyting algebras

The study of dually Stone Heyting algebras goes back to [22], while the investigations into the variety DSt of dually Stone semi-Heyting algebras were initiated in [26]. In this section we will prove that the variety $\mathbf{D S t}$ is a discriminator variety of level 1 and also present some of its properties that, besides being of interest in their own right, will be needed in the later sections. We will also consider the subvariety DStHC of DStH generated by dually Stone Heyting chains and prove that the lattice of subvarieties of $\mathbf{D S t H C}$ is an $\omega+1$-chain-a result which was implicit in [26, Section 13].

It is well-known that the identity $(x \wedge y)^{*} \approx x^{*} \vee y^{*}$ holds in Stone algebras. The following lemma is just its dual.

LEMMA 4.1. Let $\mathbf{L} \in$ DSt. Then $\mathbf{L}$ satisfies: $(x \vee y)^{\prime} \approx x^{\prime} \wedge y^{\prime}$.
See Section 2 for the definition of the condition (SC). The following theorem will be useful in the sequel.

THEOREM 4.2. Let $\mathbf{L} \in$ DSt. Then
(a) $\mathbf{L} \models x^{\prime \prime} \approx x^{\prime *}$;
(b) $\mathbf{L} \vDash(\operatorname{Lev} 1)$;
(c) If $\mathbf{L} \vDash(\mathrm{SC})$, then $\mathbf{L} \in \mathbf{J I D}_{\mathbf{1}}$.

Proof. We note that (a) is the dual of a well known property of Stone algebras. From (a) we have $\left(x \wedge x^{\prime *}\right)^{* *}=\left(x \wedge x^{\prime \prime}\right)^{* *}=x^{\prime \prime \prime *}=x^{\prime *}=$ $x^{\prime \prime}=x \wedge x^{\prime *}$, implying that (b) holds. Finally, let $\mathbf{L} \in \mathrm{DSt}$ and satisfy (SC), and let $a \in L \backslash\{1\}$. Then, by (SC) and (a), we have $a^{\prime \prime}=a \wedge a^{\prime \prime}=a \wedge a^{*}=0$, implying $a^{\prime}=1$. Then it is straightforward to verify that $\mathbf{L} \models$ (JID). Hence, (c) holds, in view of (b).

REMARK 4.3. In contrast to DSt, DPC is not, however, at level 1. For example, the algebra EIGHT with its lattice reduct, $\rightarrow$ and ' as given below, is, in fact, in the subvariety of DPC, defined by: $(x \vee y)^{\prime} \wedge\left(x^{\prime} \vee y\right)^{\prime} \wedge\left(x \vee y^{\prime}\right)^{\prime}=0$; but it fails to satisfy (Lev 1) identity.


$$
\begin{array}{r|cccccccc}
\rightarrow: & 0 & 1 & e & c & a & b & f & d \\
\hline 0 & 1 & 0 & 0 & b & b & c & 0 & 0 \\
1 & 0 & 1 & e & c & a & b & f & d \\
e & 0 & 1 & 1 & c & c & b & f & f \\
c & b & c & a & 1 & e & 0 & c & a \\
a & b & c & c & 1 & 1 & 0 & c & c \\
b & c & b & b & 0 & 0 & 1 & b & b \\
f & 0 & 1 & e & c & a & b & 1 & e \\
d & 0 & 1 & 1 & c & c & b & 1 & 1
\end{array}
$$

The following corollary is immediate from Lemma 4.1, Theorem 4.2 and [26, Corollary 8.2(a)].

COROLLARY 4.4. DSt is a discriminator variety of level 1.

Observe that Lemma 4.1 impies that DSt satisfies (BL). The following corollary is, therefore, immediate from Theorem 4.2(b) and Theorem 2.6.

COROLLARY 4.5. Let $\mathbf{L} \in \mathbf{D S t}$ with $|L| \geq 2$. Then the following are equivalent:
(1) $\mathbf{L}$ is simple,
(2) L is subdirectly irreducible,
(3) L satisfies (SC).

### 4.1. The variety DStHC.

Recall that DStHC is the variety generated by dually Stone Heyting chains. We now give an application of Corollary 4.5.
DEFINITION 4.6. For $n \in \mathbb{N}$, let $\mathbf{C}_{\mathbf{n}}^{\mathbf{d p}}$ denote the $n$-element $\mathbf{D S t H}$ chain such that
$C_{n}^{d p}=\left\{0, a_{1}, a_{2}, \ldots, a_{n-2}, 1\right\}$, where $0<a_{1}<a_{2}<\cdots<a_{n-2}<1$.
We denote by $\mathbf{V}\left(\mathbf{C}_{\mathbf{n}}^{\mathbf{d p}}\right)$ the variety generated by $\mathbf{C}_{\mathbf{n}}^{\mathbf{d p}}$. (Note that $\mathbf{C}_{\mathbf{3}}^{\mathbf{d p}}$ is the same as $\mathbf{L}_{1}^{\mathrm{dp}}$ given in [26].)

It follows from Corollary 3.3 that $\mathbf{D P C H C}=\mathbf{D S t H C}$. The following theorem was implicit in [26, Section 13].

THEOREM 4.7. The lattice of subvarieties of DStHC is the following $\omega+1$-chain:

$$
\mathbf{V}\left(\mathbf{C}_{1}^{\mathrm{dp}}\right)<\mathbf{V}\left(\mathbf{C}_{2}^{\mathrm{dp}}\right)<\cdots<\mathbf{V}\left(\mathbf{C}_{\mathbf{n}}^{\mathrm{dp}}\right)<\cdots<\text { DStHC }
$$

Proof. We claim that subdirectly irreducible algebras in DStHC are precisely the DStH-chains. For, let $\mathbf{C}^{\mathbf{d p}}$ be a DStHC-chain and let $x \in \mathbf{C}^{\mathbf{d p}}$. Since $x \leq x^{\prime}$ or $x^{\prime} \leq x$, it follows that $x=1$ or $x^{\prime}=1$, for every $x \in \mathbf{C}^{\mathbf{d p}}$, which implies that $\mathbf{C}^{\mathbf{d p}}$ satisfies (SC). On the other hand, let $\mathbf{A} \in \mathbf{D S t H C}$ satisfy (SC). Let $a \in A \backslash\{1\}$. By Theorem 4.2 (a) we have $a^{*} \leq a$; hence by (SC), we get $a^{* *}=0$, implying $a^{\prime}=1$, again by Theorem 4.2 (a). Since each DStHC-chain satisfies the identity: $(\mathrm{L})(x \rightarrow y) \vee(y \rightarrow x) \approx 1$, it follows that DStHC satisfies it too, implying that $\mathbf{A} \models(\mathrm{L})$. Hence, any two elements of $\mathbf{A}$ are comparable in $\mathbf{A}$, so $\mathbf{A}$ is a DStH-chain. Thus, $\mathbf{A} \in \mathbf{D S t H C}$ is subdirectly irreducible iff A is a DStH-chain. Now it is not hard to observe that if an identity fails in an infinite DStHC-chain, then it fails in a finite DStHC-chain. Thus DStHC is generated by finite DStH-chains. Hence, the conclusion of the theorem follows.

Note, however, that if we consider DStC-chains with semi-Heyting reducts that are not Heyting algebras, the situation gets more complicated, since the structure of the lattice of subvarieties of DStC is quite complex, as shown by the following class of examples: Let $A$ be a semi-Heyting algebra. Let $A^{e}$ be the expansion of A by adding a unary operation ' as follows:

$$
x^{\prime}=0, \text { if } x=1, \text { and } x^{\prime}=1, \text { otherwise. }
$$

Then it is clear that $A^{e}$ is a $\mathbf{D S t}$-algebra and is simple. In particular, if $A$ is a semi-Heyting-chain, then $A^{e} \in \mathbf{D S t C}$ and is simple. Furthermore, the number of semi-Heyting chains even for a small size is large; for example, there are 160 semi-Heyting chains of size 4 and, therefore, there are 160 DStC -chains of size 4 . If we denote the 2-element, non-Boolean, dually Stone semi-Heyting algebra by $\overline{2}^{\mathrm{e}}$, then it is interesting to observe that $\overline{\mathbf{2}}^{\mathrm{e}} \in \mathrm{DStC} \backslash \mathrm{DStHC}$, and $\mathbf{D S t H C}$ is only a "small" subvariety of DStC. These observations naturally suggest that the following open problem is of interest:

Problem: Investigate the structure of the lattice of subvarieties of DStC.

## 5. Subdirectly Irreducible Algebras in JID 1

Recall that the variety $\mathbf{J I D}_{1}$ is the subvriety of $\mathbf{J I D}$ defined by
(Lev 1) $\quad x \wedge x^{* *} \approx\left(x \wedge x^{\prime *}\right)^{\prime *}$.
In this section we give a somewhat concrete characterization of subdirectly irreducible (=simple) algebras in the variety $\mathrm{JID}_{1}$.

The following theorem follows immediately from Theorem 3.10 and Theorem 2.6.

THEOREM 5.1. Let $\mathbf{L} \in \mathbf{J I D}_{1}$ with $|L| \geq 2$. Then the following are equivalent:
(1) $\mathbf{L}$ is simple,
(2) $\mathbf{L}$ is subdirectly irreducible,
(3) L satisfies (SC).

We now wish to refine further the above characterization of the subdirectly irreducible algebras in $\mathrm{JID}_{1}$. In view of the above theorem, it suffices to characterize the algebras in $\mathbf{J I D}_{1}$ satisfying the condition (SC).

Unless otherwise stated, in the rest of this section we assume that $\mathrm{L} \in \mathrm{JID}_{1}$ with $|L| \geq 2$ and satisfies the simplicity condition (SC).

LEMMA 5.2. Let $a, b \in L$ such that $a^{\prime}=a$. Then

$$
a \vee b \vee b^{*}=1 .
$$

Proof. From Lemma 3.4 (4) and $a^{\prime}=a$, we have

$$
\begin{equation*}
(a \vee b) \rightarrow a=a \vee b^{*} . \tag{1}
\end{equation*}
$$

Now,

$$
\begin{array}{rlrl}
a \vee(a \vee b)^{\prime *} & =a^{\prime} \vee\left[(a \vee b)^{\prime} \rightarrow 0\right] & \\
& =\left[a^{\prime} \vee(a \vee b)^{\prime}\right] \rightarrow\left(a^{\prime} \vee 0\right), \text { by (JID) } \\
& =a^{\prime} \rightarrow a^{\prime} & \text { as } a^{\prime} \geq(a \vee b)^{\prime} \\
& =1 . &
\end{array}
$$

Thus, we have

$$
\begin{equation*}
a \vee(a \vee b)^{*}=1 \tag{2}
\end{equation*}
$$

If $a \vee b=1$, then clearly the lemma is true. So, we assume that $a \vee b \neq 1$. Then $(a \vee b) \wedge(a \vee b)^{\prime *}=0$ by (SC), and hence, we have

$$
\begin{array}{rlr}
a & =a \vee\left(b \wedge b^{*}\right) & \\
& =(a \vee b) \wedge\left(a \vee b^{*}\right) & \\
& =\left[(a \vee b) \wedge(a \vee b)^{\prime *}\right] \vee\left[(a \vee b) \wedge\left(a \vee b^{*}\right)\right] & \\
& =(a \vee b) \wedge\left[(a \vee b)^{\prime *} \vee\left(a \vee b^{*}\right)\right] & \\
& =(a \vee b) \wedge\left[(a \vee b)^{\prime *} \vee\{(a \vee b) \rightarrow a\}\right] & \\
& =(a \vee b) \wedge\left[\{(a \vee b) \rightarrow a\} \vee(a \vee b)^{*}\right] & \\
& =(a \vee b) \wedge\left[a \vee(a \vee b)^{* *}\right] & \text { by (1) } \\
& =a \vee b & \tag{2}
\end{array}
$$

Hence, $a \vee b=a$, which implies, by (1), that $a \vee b^{*}=1$. The conclusion of the lemma is now immediate.

LEMMA 5.3. Let $x \in L \backslash\{1\}$. Then $x \leq x^{\prime}$.
Proof. Since $x \neq 1$, we have $x \wedge x^{\prime *}=0$ by (SC), from which we get $\left(x^{\prime} \vee x\right) \wedge\left(x^{\prime} \vee x^{\prime *}\right)=x^{\prime}$, whence $x^{\prime} \vee x=x^{\prime}$, as $x^{\prime} \vee x^{* *}=1$ by Lemma 3.4 (3), proving the lemma.

LEMMA 5.4. Let $|L|>2$ and let $a \in L$ such that $a^{\prime}=a$. Then the height of $L$ is at most 2 .

Proof. Suppose there are $b, c \in L$ such that $0<b<c<1$. We wish to arrive at a contradiction.
From Lemma 5.3 we have $c \leq c^{\prime}$, from which it follows that

$$
\begin{equation*}
b \leq c^{\prime} . \tag{3}
\end{equation*}
$$

Claim 1: $b^{\prime}=1$.
Suppose $b^{\prime} \neq 1$. Then, by Lemma 5.3, we get $b^{\prime} \leq b^{\prime \prime} \leq b \leq c$; thus $b^{\prime} \leq c$. Next, $b \leq c$ implies $c^{\prime} \leq b^{\prime}$; and also $c \leq c^{\prime}$ from Lemma 5.3, whence $c \leq b^{\prime}$. Thus we conclude that $b^{\prime}=c$, whence $c^{\prime}=b^{\prime \prime} \leq b$, implying $c^{\prime}=b$, by (3). Then, in view of Lemma 5.3. we have $c \leq c^{\prime}=$ $b$; thus $c \leq b$, which is a contradiction, proving the claim.

From Lemma 5.2 we have $a \vee b \vee b^{*}=1$. Hence, $a^{\prime} \wedge b^{\prime} \wedge b^{* \prime}=0$ by Theorem 3.5, implying $a \wedge b^{* \prime}=0$ by Claim 1 and the hypothesis. Thus

$$
\begin{equation*}
a \wedge b^{* \prime}=0 \tag{4}
\end{equation*}
$$

Therefore, $a \vee b^{*} \geq a \vee b^{* \prime \prime}=1$ as $a^{\prime}=a$, yielding $b \leq a$. Hence, again from (4), we obtain

$$
\begin{equation*}
b \wedge b^{* \prime}=0 . \tag{5}
\end{equation*}
$$

Claim 2: $b \vee b^{*}=1$.
Suppose the claim is false. Then $b \leq b \vee b^{*} \leq\left(b \vee b^{*}\right)^{\prime}$ by Lemma 5.3, whence $b \leq b^{\prime} \wedge b^{* \prime}$, which implies $b=b \wedge b^{\prime} \wedge b^{* \prime}=0$ by the equation (5), contrary to $b>0$, proving the claim.

From Claim 2 and Theorem 3.5 we have $b^{\prime} \wedge b^{* \prime}=0$, Since $b^{\prime}=1$ by Claim 1, it follows that $b^{* \prime}=0$, whence $b^{*} \geq b^{* \prime \prime}=1$; so $b \leq b^{* *}=0$, contradicting $b>0$, proving the lemma.

LEMMA 5.5. For every $x \in \mathbf{L}, x=1$ or $x^{\prime}=1$ or $x=x^{\prime}$.
Proof. Suppose $x \in L$ such that $x \neq 1$ and $x^{\prime} \neq 1$. Then by Lemma 5.3, we have $x \leq x^{\prime}$. Also, since $x^{\prime} \neq 1$, we have $x^{\prime} \leq x^{\prime \prime} \leq x$, again by Lemma 5.3. So, $x=x^{\prime}$, proving the lemma.

LEMMA 5.6. Let $a \in L$ such that $a^{\prime}=a$. Then $a^{* \prime}=a^{*}$.
Proof. First, observe that $a \neq 0$ and $a \neq 1$, since $a=a^{\prime}$. Suppose $a^{* \prime} \neq a^{*}$. The following claims will lead to a contradiction.

Claim 1: $a^{*}=a^{* \prime \prime}$.
$a \vee a^{* \prime \prime}=a^{\prime \prime} \vee a^{* \prime \prime}=\left(a \vee a^{*}\right)^{\prime \prime}=\left[a^{\prime} \vee\left(a^{\prime} \rightarrow 0\right)\right]^{\prime \prime}=1$ by Lemma 3.4(3). Hence, $a \vee a^{* \prime \prime}=1$, implying $a^{*} \leq a^{* \prime \prime}$, and so $a^{*}=a^{* \prime \prime}$, proving the claim.

Claim 2: $a^{*}=0$.
We have, by Lemma 5.5, that $a^{* \prime}=1$ or $a^{* \prime \prime}=1$ or $a^{* \prime}=a^{* \prime \prime}$. So, by Claim 1, we get $a^{*}=a^{* \prime \prime}=0$ or $a^{*}=1\left(\right.$ as $\left.a^{*} \geq a^{* \prime \prime}\right)$ or $a^{* \prime}=a^{*}$. But, we know, by our assumption, that $a^{*} \neq a^{* \prime}$. Hence, $a^{*}=0$ or $a^{*}=1$, which clearly implies $a^{*}=0$ or $a=0$. Since we know that $a \neq 0$, the claim is proved.

Now, in view of (JID) and Claim 2, we have $a=a \vee 0=a \vee a^{*}=$ $a^{\prime} \vee(a \rightarrow 0)=\left(a^{\prime} \vee a\right) \rightarrow\left(a^{\prime} \vee 0\right)=a \rightarrow a=1$, implying $a=1$, which is a contradiction, proving the lemma.

PROPOSITION 5.7. Let $|L|>2$. Suppose there is an $a \in L$ such that $a^{\prime}=a$. Then $\mathbf{L} \in\left\{\mathbf{D}_{\mathbf{1}}, \mathbf{D}_{\mathbf{2}}, \mathbf{D}_{\mathbf{3}}\right\}$, up to isomorphism.

Proof. In view of Lemma 5.4 and $|L|>2$, the height of $L$ is exactly 2. Since the lattice reduct of $L$ is distributive, $L$ is either a 3-element chain or a 4-element Boolean lattice. We know from Lemma 5.6 that $a^{* \prime}=a^{*}$. Thus $a$ and $a^{*}$ are complementary, implying that the lattice reduct of $\mathbf{L}$ is a 4-element Boolean lattice; so $\mathbf{L} \models(B o)$, and hence $\mathbf{L} \in \mathbf{D Q B}$. Then, from Proposition 2.4 (a) it follows that $\mathbf{L} \in\left\{\mathbf{D}_{\mathbf{1}}, \mathbf{D}_{\mathbf{2}}, \mathbf{D}_{\mathbf{3}}\right\}$, up to isomorphism.

PROPOSITION 5.8. Suppose $x^{\prime} \neq x$, for every $x \in L$. Then $\mathbf{L} \in$ DSt.

Proof. Let $x \in L$. Without loss of generalty, we can assume that $x \neq 1$. Then we claim that $x^{\prime}=1$. For, assume that $x^{\prime} \neq 1$. Then, by Lemma 5.3 we get $x \leq x^{\prime}$ and $x^{\prime} \leq x^{\prime \prime}$, which implies $x=x^{\prime}$, as $x^{\prime \prime} \leq x$, contradicting the hypothesis. So, we have $x^{\prime}=1$, which implies $x^{\prime} \wedge x^{\prime \prime}=0$, Hence $\mathbf{L}$ is a dually Stone semi-Heyting algebra, completing the proof.

We are now ready to prove our main theorem of this section.
THEOREM 5.9. Let $\mathbf{L} \in$ DQD with $|L| \geq 2$. Then the following are equivalent:
(a) $\mathbf{L}$ is a subdirectly irreducible algebra in $\mathbf{J I D}_{\mathbf{1}}$,
(b) $\mathbf{L}$ is a simple algebra in $\mathbf{J I D}_{\mathbf{1}}$,
(c) $\mathbf{L} \in \mathbf{J I D}_{1}$ such that (SC) holds in $\mathbf{L}$,
(d) $\mathbf{L} \in\left\{\mathbf{D}_{\mathbf{1}}, \mathbf{D}_{\mathbf{2}}, \mathbf{D}_{\mathbf{3}}\right\}$, up to isomorphism, or $\mathbf{L} \in \mathbf{D S t}$ such that $\mathbf{L}$ satisfies (SC).

Proof. In view of Theorem 5.1, we only need to prove (c) $\Leftrightarrow$ (d). Now, suppose (d) holds. First, let us suppose $\mathbf{L} \in\left\{\mathbf{D}_{1}, \mathbf{D}_{\mathbf{2}}, \mathbf{D}_{3}\right\}$, up to isomorphism. Then it is routine to verify that $\left\{\mathbf{D}_{\mathbf{1}}, \mathbf{D}_{\mathbf{2}}, \mathbf{D}_{3}\right\} \subseteq \mathbf{J I D}_{1}$ and $\left\{\mathbf{D}_{\mathbf{1}}, \mathbf{D}_{\mathbf{2}}, \mathbf{D}_{\mathbf{3}}\right\}$ satisfies (SC), implying (c). Next, suppose $\mathbf{L} \in \mathbf{D S t}$ such that $\mathbf{L}$ satisfies (SC). Then $\mathbf{L} \in \mathbf{J I D}_{\mathbf{1}}$, in view of Theorem 4.2(c), implying that (c) holds. Thus (d) $\Rightarrow$ (c). To prove the converse, suppose (c) holds. We consider two cases. First, suppose there is an $a \in L$ such that $a^{\prime}=a$. Then, by Proposition 5.7, $\mathbf{L} \in\left\{\mathbf{D}_{\mathbf{1}}, \mathbf{D}_{\mathbf{2}}, \mathbf{D}_{\mathbf{3}}\right\}$, up to isomorphism, implying (d).

Next, suppose there is no element $a \in L$ such that $a^{\prime}=a$. Hence, $\mathbf{L}$ satisfies:

$$
\begin{equation*}
\text { For every } x \in L, x^{\prime} \neq x \tag{6}
\end{equation*}
$$

Then, using Proposition 5.8, we obtain that $\mathbf{L}$ is dually Stone, which, together with the hypothesis, leads us to conclude $(\mathrm{c}) \Rightarrow(\mathrm{d})$.

We have the following important consequence of Theorem 5.9.
COROLLARY 5.10. $\mathrm{JID}_{1}=\mathrm{DSt} \vee \mathrm{V}\left(\mathrm{D}_{1}, \mathrm{D}_{2}, \mathrm{D}_{3}\right)$.
Recall that JIDH is the subvariety of JID defined by the identity: $(x \wedge y) \rightarrow x \approx 1$, and DStH is the variety of dually Stone Heyting algebras. Now, we focus on the subvariety $\mathbf{J I D H}_{1}$ of JIDH. Note that the variety of Boolean algebras is the only atom in the lattice of subvarieties of $\mathbf{J I D H}_{\mathbf{1}}$. For $\mathbf{V}$ a subvariety of $\mathbf{J I D H}_{\mathbf{1}}$, let $\mathcal{L}(\mathbf{V})$ and $\mathcal{L}^{+}(\mathbf{V})$ denote, respectively, the lattice of subvarieties of $\mathbf{V}$ and the
lattice of nontrivial subvarieties of $\mathbf{V}$. Let $\mathbf{1} \oplus \mathbf{L}$ denote the ordinal sum of the trivial lattice $\mathbf{1}$ and a lattice $\mathbf{L}$.

Restricting the semi-Heyting reduct in the above corollary to Heyting algebras, we obtain the following interesting corollary, where $\mathbf{2}$ denotes a 2-element lattice.

COROLLARY 5.11. We have
(1) $\mathbf{J I D H}_{\mathbf{1}}=\mathbf{D S t H} \vee \mathbf{V}\left(\mathbf{D}_{\mathbf{2}}\right)$,
(2) $\mathcal{L}\left(\mathbf{J I D H}_{1}\right) \cong \mathbf{1} \oplus\left(\mathcal{L}^{+}(\mathbf{D S t H}) \times \mathbf{2}\right)$.

The preceding corollary leads to the following open problem.
PROBLEM: Investigate the structure of $\mathcal{L}^{+}(\mathrm{DStH})$.

## 6. JI-distributive, dually quasi-De Morgan, linear Semi-Heyting Algebras

In this section we focus on the linear identity:
(L) $\quad(x \rightarrow y) \vee(y \rightarrow x) \approx 1$.

Let DQDL [JIDL] denote the subvariety of DQD [JID] defined by (L), and let JIDLH denote the subvariety of JIDL consisting of JI-distributive, dually quasi-De Morgan, linear Heyting algebras.

The following result is needed later in this section. Part (a) of it is proved in [26, Lemma 12.1(f)], and (b) follows immediately from (a).

PROPOSITION 6.1. [26, Lemma 12.1(f)] Let $\mathbf{L}$ be a linear semiHeyting algebra (i.e., $\mathbf{L} \models(\mathrm{L})$ ). Then
(a) $\mathbf{L} \models(\mathrm{H})$,
(b) $\mathrm{JIDL}=\mathrm{JIDLH}$.

LEMMA 6.2. Let $\mathbf{L} \in \mathrm{DQDL}$ and let $x, y \in L$. Then
(a) $(x \rightarrow y) \vee(y \rightarrow x)^{\prime \prime}=1$,
(b) $x \leq y \vee(y \rightarrow x)^{\prime \prime}$.

Proof. $(x \rightarrow y) \vee(y \rightarrow x)^{\prime \prime} \geq(x \rightarrow y)^{\prime \prime} \vee(y \rightarrow x)^{\prime \prime}=[(x \rightarrow y) \vee(y \rightarrow$ $x)]^{\prime \prime}=1$ by (L), proving (a). Using (a), we get $x \wedge\left[y \vee(y \rightarrow x)^{\prime \prime}\right]=$ $(x \wedge y) \vee\left[x \wedge(y \rightarrow x)^{\prime \prime}\right]=\left[x \wedge(x \rightarrow y) \vee\left[x \wedge(y \rightarrow x)^{\prime \prime}\right]=x \wedge[(x \rightarrow\right.$ $\left.y) \vee(y \rightarrow x)^{\prime \prime}\right]=x \wedge 1=x$, implying (b).

Note that the algebra SIX described earlier in Section 3 is actually an algebra in JIDL. Hence JIDL does not satisfy (Lev 1); but JIDL is at level 2, in view of Theorem 3.10.

In this section, our goal is to present, as an application of Theorem 5.9, an explicit description of subdirectly irreducible ( $=$ simple) algebras in the variety $\mathrm{JIDL}_{1}$ of JI-distributive, dually quasi-De Morgan, linear semi-Heyting algebras of level 1.

Recall that $\mathrm{DPCC}=\mathrm{DStC}$ and $\mathrm{DPCHC}=\mathrm{DStHC}$. So, we use these names interchangeably.

LEMMA 6.3. DPCC $\models$ (Lev1).
Proof. Let L be a DPC-chain and let $x \in L$. Since $x, x^{\prime}$ are comparable, we have $x \vee x^{\prime}=x$ or $x \vee x^{\prime}=x^{\prime}$, implying $x=1$ or $x^{\prime}=1$, as $x \vee x^{\prime}=1$. Then it is easy to see that (Lev 1) holds in $\mathbf{L}$, and hence in DPCC.

## PROPOSITION 6.4. DPCHC $\vee \mathrm{V}\left(\mathrm{D}_{2}\right) \subseteq \mathrm{JIDL}_{1}$.

Proof. It follows from Lemma 3.2, and Lemma 6.3 that DPCHC satisfies (JID) and (Lev 1), and it is easy to see that DPCHC $\models(\mathrm{L})$. Also, it is routine to verify that (JID), (L) and (Lev 1) hold in $\mathbf{D}_{\mathbf{2}}$.

Our goal in this section is to prove that, in fact, the equality holds in the statement of the above Proposition.

Unless otherwise stated, in the rest of this section we assume that $\mathrm{L} \in \mathrm{JIDL}_{1}$ with $|L|>2$ and satisfies (SC).

LEMMA 6.5. Let $x, y \in L$ such that $x \vee y \neq 1$. Then, $x \leq y^{\prime}$.
Proof. Let $x \vee y \neq 1$. Since $y^{\prime} \vee(x \vee y)^{\prime *} \geq y^{\prime} \vee y^{\prime *}=1$ by Lemma 3.4 (3), we get, using (SC), that $x=x \wedge(x \vee y) \wedge\left[y^{\prime} \vee(x \vee y)^{* *}\right]=$ $x \wedge\left[\left\{(x \vee y) \wedge y^{\prime}\right\} \vee\left\{(x \vee y) \wedge(x \vee y)^{\prime *}\right\}\right]=x \wedge(x \vee y) \wedge y^{\prime}=x \wedge y^{\prime}$, whence $x \leq y^{\prime}$.

LEMMA 6.6. Let $a, b \in L$ such that $a^{\prime} \neq a, a \neq 1$, and $a \not \leq b$. Then $(a \rightarrow b)^{\prime \prime}=0$.
Proof. First, we claim that $a \not \leq(a \rightarrow b)^{\prime \prime}$. For, suppose $a \leq(a \rightarrow b)^{\prime \prime}$; then $a=a \wedge(a \rightarrow b)^{\prime \prime} \leq b$ by Lemma 2.2 (vii), implying $a \leq b$, which is a contradiction to the hypothesis $a \not \leq b$. Hence $a \not \leq(a \rightarrow b)^{\prime \prime}$. Then $a \vee(a \rightarrow b)^{\prime}=1$ by (the contrapositive of) Lemma 6.5, whence $a^{\prime \prime} \vee(a \rightarrow b)^{\prime \prime \prime}=1$. Since $a \neq 1$ and $a^{\prime} \neq a$ by hypothesis, we get $a^{\prime}=1$ by Lemma 5.5, whence $a^{\prime \prime}=0$. Then we conclude that $(a \rightarrow b)^{\prime \prime}=0$, proving the lemma.

We are now ready to give an explicit description of subdirectly irreducible (=simple) algebras in $\mathbf{J I D L}_{\mathbf{1}}$.

THEOREM 6.7. Let $\mathbf{L} \in \mathrm{DQD}_{1}$ with $|L|>2$. Then the following are equivalent:
(1) $\mathbf{L}$ is a subdirectly irreducible algebra in $\mathbf{J I D L}_{\mathbf{1}}$,
(2) $\mathbf{L}$ is a simple algebra in $\mathbf{J I D L}_{\mathbf{1}}$,
(3) $\mathbf{L} \in \mathbf{J I D L}_{\mathbf{1}}$ such that (SC) holds in $\mathbf{L}$,
(4) $\mathbf{L} \cong \mathbf{D}_{\mathbf{2}}$, or $\mathbf{L}$ is a $\mathbf{D S t H}$-chain.

Proof. (1) $\Leftrightarrow(2) \Leftrightarrow$ (3) follow from Theorem 5.9. So we need to prove (3) $\Rightarrow(4) \Rightarrow(3)$. Suppose (3) holds. Then, by Theorem 5.9, either $\mathbf{L} \in\left\{\mathbf{D}_{\mathbf{1}}, \mathbf{D}_{\mathbf{2}}, \mathbf{D}_{\mathbf{3}}\right\}$, or $\mathbf{L} \in \mathbf{D S t}$ and satisfies (SC). In the former case, since $\mathbf{L} \models(\mathrm{L})$, it follows from Proposition 6.1 that $\mathbf{L} \models(\mathrm{H})$. Hence $\mathbf{L} \cong \mathbf{D}_{\mathbf{2}}$. Next, we assume the latter case. So, $\mathbf{L} \in \mathbf{D S t}$ and satisfies (SC). Since $\mathbf{L} \models(\mathrm{L})$ by hypothesis, we get, by Proposition 6.1, that $\mathbf{L} \in \mathbf{D S t H}$. So, we need only prove that $\mathbf{L}$ is a chain. Let $a, b \in L \backslash\{1\}$ such that $a \not \leq b$. Then, from Lemma 6.2(b), we have that $b \leq a \vee(a \rightarrow b)^{\prime \prime}$. Also, since $\mathbf{L} \models(\mathrm{DSt})$, it is clear that $a^{\prime} \neq a$. Hence, by Lemma 6.6, we get $(a \rightarrow b)^{\prime \prime}=0$, implying $b \leq a$. Thus, the lattice reduct of $\mathbf{L}$ is a chain, and so, $(3) \Rightarrow(4)$. Finally, assume (4) holds. First, if $\mathbf{L} \cong \mathbf{D}_{\mathbf{2}}$, then it is routine to verify that (3) holds. Next, suppose $\mathbf{L}$ is a $\mathbf{D S t H}$-chain and let $x \in L$. Then, $x^{\prime} \leq x^{\prime \prime}$ or $x^{\prime \prime} \leq x^{\prime}$, implying $x^{\prime} \wedge x^{\prime \prime}=x^{\prime}$ or $x^{\prime} \wedge x^{\prime \prime}=x^{\prime \prime}$. Hence, by ( DSt ), we get $x^{\prime}=0$ or $x^{\prime}=1$, from which it is easy to see that $\mathbf{L}$ satisfies (SC). So, from Theorem 4.2 (c), we conclude that $\mathbf{L} \in \mathbf{J I D}_{\mathbf{1}}$. Also, it is well known that Heyting chains satisfy (L). Thus, $\mathbf{L} \in \mathbf{J I D L}_{\mathbf{1}}$ and $\mathbf{L}$ satisfies (SC), implying (3).

The following corollary is immediate from Theorem 6.7.

## COROLLARY 6.8. $\mathrm{JIDL}_{1}=\mathrm{DStHC} \vee \mathrm{V}\left(\mathrm{D}_{2}\right)$.

We would like to mention here that the attempt to solve the problem of axiomatization of $\mathbf{D S t H C} \vee \mathrm{V}\left(\mathbf{D}_{\mathbf{2}}\right)$ led to the results of this paper, with Corollary 6.8 yielding a solution to that problem.

We conclude this section with an axiomatization of DStHC.

## THEOREM 6.9. DStHC $=$ DStL .

Proof. We know from the proof of Theorem 4.7 that the subdirectly irreducible algebras in DStHC are precisely the DStH-chains. So, to complete the proof, it suffices to prove that the subdirectly irreducible algebras in DStL are precisely the DStH-chains. For this, first note that from Proposition 6.1 we have that a linear semi-Heyting algebra satisfies (L) and hence DStL $\models(H)$, implying DStL $\subseteq$ DStH. Now, let $\mathbf{L}$ be a subdirectly irreducible ( $=$ simple) algebra in DStL. We wish to show that $\mathbf{L}$ is a Heyting chain. Let $a, b \in L$ be arbitrary. By

Corollary 4.5, $\mathbf{L}$ satisfies (SC); and $\mathbf{L} \models(\mathrm{H})$, as $\mathbf{D S t L} \models(\mathrm{H})$. Hence, by Lemma 6.2 (a), we have

$$
\begin{equation*}
\mathbf{L} \models(x \rightarrow y) \vee(y \rightarrow x)^{\prime \prime}=1 . \tag{7}
\end{equation*}
$$

Suppose that $a \neq 1$. Then from (SC) we have $a \wedge a^{\prime *}=0$, whence $a \wedge a^{\prime \prime}=0$ by Lemma 4.2 (a), implying $a^{\prime}=1$. Thus we have proved

For every $x \in L, x=1$ or $x^{\prime}=1$.
Hence, by (8), we get $(a \rightarrow b)^{\prime}=1$ or $(a \rightarrow b)^{\prime \prime}=1$, implying $(a \rightarrow$ $b)^{\prime \prime}=0$ or $(a \rightarrow b)=1$. So, by (7), we have $b \rightarrow a=1$ or $a \rightarrow b=1$, implying $b \leq a$ or $a \leq b$, as $\mathbf{L} \models(\mathrm{H})$. Thus $\mathbf{L}$ is a DStH-chain, completing the proof.

## 7. More Consequences of Theorem 6.7

In this section we present some more consequences of Theorem 6.7.
As mentioned earlier, the axiomatizations of the variety DPCHC (= DStHC) and all of its subvarieties were given in [26].

The following corollary is immediate from Corollary 6.8 and Theorem 4.7.

COROLLARY 7.1.
(1) $\mathcal{L}\left(\mathrm{JIDL}_{1}\right) \cong \mathbf{1} \oplus[(\omega+\mathbf{1}) \times \mathbf{2}]$.
(2) $\mathrm{JIDL}_{1}$ and $\mathbf{D S t H C}$ are the only two elements of infinite height in the lattice $\mathcal{L}\left(\mathbf{J I D L}_{1}\right)$.
(3) $\mathbf{V} \in \mathcal{L}^{+}\left(\mathbf{J I D L}_{\mathbf{1}}\right)$ is of finite height iff $\mathbf{V}$ is either $\mathbf{V}\left(\mathbf{D}_{\mathbf{2}}\right)$, or $\mathbf{V}\left(\mathbf{C}_{\mathbf{n}}^{\mathbf{d p}}\right)$ for some $n \in \mathbb{N} \backslash\{1\}$, or $\mathbf{V}\left(\mathbf{C}_{\mathbf{m}}^{\mathbf{d p}}\right) \vee \mathbf{V}\left(\mathbf{D}_{\mathbf{2}}\right)$ for some $m \in \mathbb{N} \backslash\{1\}$.
In Corollaries 7.2-7.5, we give equational bases to all subvarieties of $\mathrm{JIDL}_{1}$.

COROLLARY 7.2. The variety $\operatorname{DStHC}$ is defined, modulo $\mathrm{JIDL}_{1}$, by

$$
x \vee x^{\prime} \approx 1
$$

Proof. Observe that $\mathbf{D S t H C} \models x \vee x^{\prime} \approx 1$, but $\mathbf{V}\left(\mathbf{D}_{\mathbf{2}}\right) \not \vDash x \vee x^{\prime} \approx 1$, and then apply Theorem 6.7.

The variety $\mathbf{V}\left(\mathbf{D}_{2}\right)$ was axiomatized in [26]. Here is a new one.
COROLLARY 7.3. The variety $\mathbf{V}\left(\mathbf{D}_{2}\right)$ is defined, modulo $\mathbf{J I D L}_{1}$, by

$$
x^{\prime \prime} \approx x
$$

Proof. Observe that $\mathbf{D S t H C} \not \vDash x^{\prime \prime} \approx x$, but $\mathbf{V}\left(\mathbf{D}_{\mathbf{2}}\right) \neq x^{\prime \prime} \approx x$, and then use Theorem 6.7.

COROLLARY 7.4. Let $n \geq 2$. The variety $\mathbf{V}\left(\mathbf{C}_{\mathbf{n}}^{\mathbf{d p}}\right) \vee \mathbf{V}\left(\mathbf{D}_{\mathbf{2}}\right)$ is defined, modulo $\mathbf{J I D L}_{\mathbf{1}}$, by
$\left(\mathrm{C}_{n}\right) x_{1} \vee x_{2} \vee \cdots \vee x_{n} \vee\left(x_{1} \rightarrow x_{2}\right) \vee\left(x_{2} \rightarrow x_{3}\right) \vee \cdots \vee\left(x_{n-1} \rightarrow x_{n}\right)=1$.
Proof. We now prove that $\mathbf{C}_{\mathbf{n}}^{\mathbf{d p}} \models\left(\mathrm{C}_{\mathbf{n}}\right)$. Let $\left\langle c_{1}, c_{2}, \ldots, c_{n}\right\rangle \in C_{n}^{d p}$ be an arbitrary assignment in $C_{n}^{d p}$ for the variables $x_{i}$ such that $c_{i}$ is the value of $x_{i}$, for $i=1, \cdots, n$. If $c_{i} \leq c_{i+1}$ for some $i$, then $c_{i} \rightarrow c_{i+1}=1$, as $\mathbf{C}_{\mathbf{n}}^{\mathbf{d p}}$ has a Heyting algebra reduct, and hence, the identity holds in $\mathbf{C}_{\mathbf{n}}^{\mathbf{d p}}$. So, we assume that $c_{i}>c_{i+1}$, for $i=1,2, \cdots, n$. Then, $c_{1}=1$ since $\left|C_{n}^{d p}\right|=n$, implying that $\left(\mathrm{C}_{n}\right)$ holds in $\mathbf{C}_{\mathbf{n}}^{\mathrm{dp}}$. Also, it is routine to check that that $\mathbf{D}_{2} \models\left(\mathrm{C}_{2}\right)$ and $\left(\mathrm{C}_{i}\right)$ implies $\left(\mathrm{C}_{i+1}\right)$, for $i=2, \cdots, n-1$. So, $\mathbf{D}_{2} \models\left(\mathbf{C}_{\mathbf{n}}\right)$, implying that $\mathbf{V}\left(\mathbf{C}_{\mathbf{n}}^{\mathbf{d p}}\right) \vee \mathbf{V}\left(\mathbf{D}_{\mathbf{2}}\right) \models\left(\mathrm{C}_{\mathbf{n}}\right)$.

Next, suppose that $\mathbf{V}$ is the subvariety of $\mathbf{J I D L}_{1}$ satisfying $\left(\mathrm{C}_{n}\right)$. Then, by Corollary $3.11, \mathbf{V}$ is a discriminator variety. Let $\mathbf{L}$ be a simple ( $=$ subdirectly irreducible) algebra in V. Then, it follows from Corollary 6.8 (or Theorem 6.7) that $\mathbf{L}$ is a $\mathbf{D S t H}$-chain or $\mathbf{L} \cong \mathbf{D}_{\mathbf{2}}$. Suppose that $\mathbf{L}$ is a $\mathbf{D S t H}$-chain. Assume, if possible, $|L|>n$. Then, there exist $b_{1}, b_{2} \cdots, b_{n-1} \in L$ such that $0<b_{1}<\cdots,<$ $b_{n-1}<1$. Since $\mathbf{L} \models\left(\mathrm{C}_{n}\right)$, we can assign $\left\langle b_{n-1}, b_{n-2}, \cdots, b_{1}, 0\right\rangle$ for $\left\langle x_{1}, x_{2}, \cdots, x_{n-1}, x_{n}\right\rangle$. Then, $b_{n-1} \vee\left(b_{n-1} \rightarrow b_{n-2}\right) \vee \cdots, \vee\left(b_{1} \rightarrow 0\right)=1$, yielding $b_{n-1} \vee b_{n-2} \vee \cdots \vee b_{1} \vee 0=1$, implying that $b_{n-1}=1$, which is a contradiction. Thus we have $|L| \leq n$, from which it follows that $\mathbf{V} \subseteq \mathbf{V}\left(\mathbf{C}_{\mathbf{n}}^{\mathrm{dp}}\right) \vee \mathbf{V}\left(\mathbf{D}_{\mathbf{2}}\right)$, completing the proof.
COROLLARY 7.5. The variety $\mathbf{V}\left(\mathbf{C}_{\mathbf{n}}^{\mathbf{d p}}\right)$ is defined, modulo $\mathbf{J I D L}_{\mathbf{1}}$, by
(1) $x \vee x^{\prime} \approx 1$,
(2) $x_{1} \vee x_{2} \vee \cdots \vee x_{n} \vee\left(x_{1} \rightarrow x_{2}\right) \vee\left(x_{2} \rightarrow x_{3}\right) \vee \cdots \vee\left(x_{n-1} \rightarrow x_{n}\right)=1$.

For a different base for $\mathbf{V}\left(\mathbf{C}_{\mathbf{n}}^{\mathbf{d p}}\right)$, see [26]. Regularity was studied in [26], [27], [28] and [29]. Here is another use of it.
COROLLARY 7.6. The variety $\mathbf{V}\left(\mathbf{C}_{\mathbf{3}}^{\mathrm{dp}}\right) \vee \mathbf{V}\left(\mathbf{D}_{\mathbf{2}}\right)$ is defined, modulo $\mathrm{JIDL}_{1}$, by
$x \wedge x^{+} \leq y \vee y^{*}$ (Regularity), where $x^{+}:=x^{\prime * \prime}$.
It is also defined, modulo $\mathbf{J I D L}_{\mathbf{1}}$, by

$$
x \wedge x^{\prime} \leq y \vee y^{*}
$$

The variety $\mathbf{V}\left(\mathbf{C}_{3}^{\mathrm{dp}}\right)$ is axiomatized in [26]. Here is another axiomatization for it.

COROLLARY 7.7. The variety $\mathbf{V}\left(\mathbf{C}_{3}^{\mathrm{dp}}\right)$ is defined, modulo $\mathbf{J I D L}_{\mathbf{1}}$, by
(1) $x \wedge x^{+} \leq y \vee y^{*}$ (Regularity),
(2) $x^{\prime}=x^{+}$.
7.1. Amalgamation Property. We now examine the Amalgamation Property for subvarieties of the variety DStHC. For this purpose, we need the following lemma whose proof is straightforward.

We use " $\leq$ " to abbreviate "is a subalgebra of" in the next lemma.
Recall from Theorem 4.7 (see also [26]) that the proper, nontrivial subvarieties of DStHC are precisely the subvarieties of the form $\mathbf{V}\left(\mathbf{C}_{\mathbf{n}}^{\mathbf{d p}}\right)$, for $n \in \mathbb{N}$.
LEMMA 7.8. Let $m, n \in \mathbb{N}$. Then

$$
\mathbf{C}_{\mathbf{m}}^{\mathbf{d p}} \leq \mathbf{C}_{\mathbf{n}}^{\mathbf{d p}}, \text { for } m \leq n
$$

COROLLARY 7.9. Every subvariety of DStHC has Amalgamation Property.

Proof. It follows from Corollary 4.4 that DStHC is a discriminator variety; and hence has CEP. Also, from Theorem 6.7 we obtain that every subalgebra of each subdirectly irreducible ( $=$ simple) algebra in DStHC is subdirectly irreducible. Let V be a subvariety of DStHC. Then, using a result from [11] that we need only consider an amalgam $(\mathbf{A}: \mathbf{B}, \mathbf{C})$, where $\mathbf{A}, \mathbf{B}, \mathbf{C}$ are simple in $\mathbf{V}$ and $\mathbf{A}$ a subalgebra of $\mathbf{B}$ and $\mathbf{C}$. First, suppose $\mathbf{V}=\mathbf{V}\left(\mathbf{C}_{\mathbf{n}}^{\mathbf{d p}}\right)$ for some $n$. Then $\mathbf{B}$ and $\mathbf{C}$ are DStHC-chains. Then, in view of the preceding lemma, it is clear that the amalgam $(\mathbf{A}: \mathbf{B}, \mathbf{C})$ can be amalgamated in $\mathbf{V}$. Next, suppose $\mathbf{V}=\mathbf{D S t H C}$, then it is clear that the amalgamation can be achieved as in the previous case.

We conclude this section with the following remark: Since every subvariety V of DStHC has Congruence Extension Property and Amalgamation Property, it follows from Banachewski [6] that all subvarieties of DStHC have enough injectives (see [6] for the definition of this notion).

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# OPEN MAPPING THEOREMS WITH FINITE FIBRES FOR $C$-SPACES 

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#### Abstract

In this paper we study theorems for $C$-spaces and finite $C$-spaces on dimensionraising open mappings and dimension-lowering open mappings with finite fibres.


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## 1 Introduction

In this paper we assume that all spaces are normal and all mappings are continuous.
A space $X$ is $A$-weakly infinite-dimensional or Alexandroff weakly infinite-dimensional if for every collection $\left\{\left(A_{i}, B_{i}\right): i<\omega\right\}$ of pairs of disjoint closed subsets of $X$ there exists a collection $\left\{L_{i}: i<\omega\right\}$ of closed subsets of $X$ such that $L_{i}$ is a partition in $X$ between $A_{i}$ and $B_{i}$ for every $i<\omega$, and $\bigcap_{i<\omega} L_{i}=\emptyset$.

A space $X$ is a $C$-space [1] if for every countable collection $\left\{\mathcal{G}_{i}: i<\omega\right\}$ of open covers of $X$ there exists a countable collection $\left\{\mathcal{H}_{i}: i<\omega\right\}$ of collections of pairwise disjoint open subsets of $X$ such that $\mathcal{H}_{i}$ is a refinement of $\mathcal{G}_{i}$ for every $i<\omega$ and $\bigcup_{i<\omega} \bigcup\left\{H: H \in \mathcal{H}_{i}\right\}=X$.

It is easily seen that every $C$-space is A-weakly infinite-dimensional. However, it is still unknown whether every compact A-weakly infinite-dimensional metrizable space is a $C$-space.

A space $X$ is S-weakly infinite-dimensional or Smirnov weakly infinite-dimensional if for every collection $\left\{\left(A_{i}, B_{i}\right): i<\omega\right\}$ of pairs of disjoint closed subsets of $X$ there exists a collection $\left\{L_{i}: i<\omega\right\}$ of closed subsets of $X$ such that $L_{i}$ is a partition in $X$ between $A_{i}$ and $B_{i}$ for every $i<\omega$, and $\bigcap_{i \leq n} L_{i}=\emptyset$ for some $n<\omega$. It directly follows from the definition that every S-weakly infinite dimensional space is A-weakly infinite dimensional, and every compact A-weakly infinite dimensional space is S-weakly infinite dimensional.

A space $X$ is a finite $C$-space [2] if for every collection $\left\{\mathcal{G}_{i}: i<\omega\right\}$ of finite open covers of $X$ there exists a collection $\left\{\mathcal{H}_{i}: i<\omega\right\}$ of collections of pairwise disjoint open subsets of $X$ such that $\mathcal{H}_{i}$ is a refinement of $\mathcal{G}_{i}$ for every $i<\omega$ and $\bigcup_{i \leq n} \bigcup\left\{H: H \in \mathcal{H}_{i}\right\}=X$ for some $n<\omega$. It is well-known [2] that every finite $C$-space is S -weakly infinite-dimensional. There exists a $C$-space which is not a finite $C$-space (see [1, Example 2.15]). However, every compact $C$-space is a finite $C$-space.

For paracompact spaces, Gutev and Valov [5] proved the countable sum theorem for $C$-spaces. For countably paracompact and collectionwise normal spaces, the author proved the countable sum theorem for $C$-spaces (cf. [6, Corollary 3.2]). Addis and Gresham [1] proved that every finite-dimensional, paracompact space is a $C$-space. By the same proof, we can show that every finite-dimensional space is a finite $C$-spaces. The following two Lemmas will play a important role in the proof of our main theorems.

Lemma A If there exists a closed subset $K$ of a countably paracompact collectionwise normal space $X$ satisfying the following conditions (1) and (2), then $X$ is a $C$-space.
(1) $K$ is a $C$-space,
(2) for every closed subset $F$ of $X$ with $F \cap K=\emptyset, F$ is a $C$-space.

Proof. Let $\left\{\mathcal{G}_{i}: i<\omega\right\}$, where $\mathcal{G}_{i}=\left\{G_{\lambda}: \lambda \in \Lambda_{i}\right\}$, be a collection of open covers of $X$. Since $K$ is a countably paracompact $C$-space, by [6, Lemma 2.1], there exists a collection $\left\{\mathcal{U}_{2 i}: i<\omega\right\}$, where $\mathcal{U}_{2 i}=\left\{U_{\lambda}: \lambda \in \Lambda_{2 i}\right\}$, of discrete collections of open
subsets of $K$ such that $U_{\lambda} \subset G_{\lambda} \cap K$ and $\bigcup_{i<\omega} \bigcup\left\{U_{\lambda}: \lambda \in \Lambda_{2 i}\right\}=K$. Since $X$ is collectionwise normal, by [6, Lemma 2.2], there exists a collection $\left\{\mathcal{H}_{2 i}: i<\omega\right\}$, where $\mathcal{H}_{2 i}=\left\{H_{\lambda}: \lambda \in \Lambda_{2 i}\right\}$, of discrete collections of open subsets of $X$ such that $H_{\lambda} \cap K=U_{\lambda}$ and $H_{\lambda} \subset G_{\lambda}$. Let us set $F=X-\bigcup_{i<\omega} \bigcup\left\{H: H \in \mathcal{H}_{2 i}\right\}$. Similarly there exsits a collection $\left\{\mathcal{H}_{2 i+1}: i\langle\omega\}\right.$ of discrete collections of open subsets of $X$ such that $\mathcal{H}_{2 i+1}$ is a refinement of $\mathcal{G}_{2 i+1}$ for every $i<\omega$ and $\bigcup_{i<\omega} \bigcup\left\{H: H \in \mathcal{H}_{2 i+1}\right\} \supset F$. We get the required collection $\left\{\mathcal{H}_{i}: i<\omega\right\}$.

Lemma B If there exists a closed subset $K$ of a space $X$ satisfying the following conditions (1) and (2), then $X$ is a finite $C$-space.
(1) $K$ is a finite $C$-space,
(2) for every closed subset $F$ of $X$ with $F \cap K=\emptyset, F$ is a finite $C$-space.

Proof. Let $\left\{\mathcal{G}_{i}: i<\omega\right\}$, where $\mathcal{G}_{i}=\left\{G_{\lambda}: \lambda \in \Lambda_{i}\right\}$, be a collection of finite open covers of $X$. Since $K$ is a finite $C$-space, there exists a collection $\left\{\mathcal{U}_{2 i}: i<\omega\right\}$, where $\mathcal{U}_{2 i}=\left\{U_{\lambda}: \lambda \in \Lambda_{2 i}\right\}$, of finite collections of pairwise disjoint open subsets of $K$ such that $U_{\lambda} \subset G_{\lambda} \cap K$ and $\bigcup_{i=1}^{n} \bigcup\left\{U_{\lambda}: \lambda \in \Lambda_{2 i}\right\}=K$ for some $n<\omega$. Since $K$ is normal, there exists $\left\{\mathcal{F}_{2 i}: i \leq n\right\}$, where $\mathcal{F}_{2 i}=\left\{F_{\lambda}: \lambda \in \Lambda_{2 i}\right\}$, of collections of closed subsets of $K$ such that $F_{\lambda} \subset U_{\lambda}$ and $\bigcup_{i=1}^{n} \bigcup\left\{F_{\lambda}: \lambda \in \Lambda_{2 i}\right\}=K$. There exists a collection $\left\{\mathcal{H}_{2 i}: i<\omega\right\}$, where $\mathcal{H}_{2 i}=\left\{H_{\lambda}: \lambda \in \Lambda_{2 i}\right\}$, of finite collections of pairwise disjoint open subsets of $X$ for every $i<\omega$ such that $F_{\lambda} \subset H_{\lambda} \subset G_{\lambda}$ and $\bigcup_{i=1}^{n} \bigcup\left\{H_{\lambda}: \lambda \in \Lambda_{2 i}\right\} \supset K$. For every $i>n$ we let $\mathcal{H}_{2 i}=\{\emptyset\}$. Let us set $F=X-\bigcup_{i \leq n} \bigcup\left\{H: H \in \mathcal{H}_{2 i}\right\}$. For a space $F$ repeating above procedure we obtain the required collection $\left\{\mathcal{H}_{i}: i<\omega\right\}$.

## 2 Dimension-raising mappings

Polkowski [8] proved the following theorem.
Theorem [8]. If $f: X \longrightarrow Y$ is an open mapping of an $A$-weakly infinite-dimensional space $X$ onto a countably paracompact space $Y$ such that $\left|f^{-1}(y)\right|<\omega$ for every $y \in Y$, then $Y$ is $A$-weakly infinite-dimensional.

We shall prove the following theorem. This is an analogy of the above Polkowski's theorem.
2.1. Theorem If $f: X \longrightarrow Y$ is an open mapping of a $C$-space $X$ onto a countably paracompact and collectionwise normal $Y$ such that $\left|f^{-1}(y)\right|<\omega$ for every $y \in Y$, then $Y$ is a $C$-space.

To prove Theorem 2.1 we need the following theorem and lemma.
2.2. Theorem(cf.[4, Lemma 6.7]) If $f: X \longrightarrow Y$ is a closed mapping of a countably paracompact $C$-space $X$ onto a space $Y$ and there exists an integer $k \geq 1$ such that $\left|f^{-1}(y)\right| \leq k$ for every $y \in Y$, then $Y$ is a $C$-space.
2.3. Lemma([3, Lemma 6.3.12]) If all fibres of an open mapping $f: X \longrightarrow Y$ defined on a space $X$ are finite and have the same cardinality, then $f$ is closed.
2.4 Proof of theorem 2.1. Let $K_{j}=\left\{y \in Y:\left|f^{-1}(y)\right|=j\right\}$ for every $j \in \mathbb{N}$. It is easy to see that the union $\bigcup_{j \leq i} K_{j}$ is closed in $Y$ for every $i \in \mathbb{N}$. Inductively, we show that the union $\bigcup_{j \leq i} K_{j}$ is a $C$-space for every $i \in \mathbb{N}$. To this end, it suffices to show that every closed subspace $Z$ of $Y$ contained in $K_{i}$ is a $C$-space, cf. Lemma A. By Lemma 2.3, the restriction $\left.f\right|_{f^{-1}(Z)}: f^{-1}(Z) \longrightarrow Z$ is perfect. As the inverse image of a countably paracompact space under a perfect mapping is countably paracompact, then $f^{-1}(Z)$ is countably paracompact. By Theorem $2.2, Z$ is a $C$-space. Thus the union $\bigcup_{j \leq i} K_{j}$ is a closed $C$-space for every $i \in \mathbb{N}$. By countable sum theorem, $Y$ is a $C$-space.

The following theorem is a counterpart for finite $C$-spaces of Polkowski's result.
2.5. Theorem If $f: X \longrightarrow Y$ is an open mapping of a weakly paracompact finite $C$-space $X$ onto a space $Y$ such that $\left|f^{-1}(y)\right|<\omega$ for every $y \in Y$, then $Y$ is a finite $C$-space.

To prove Theorem 2.5 we need the following theorem and lemma.
2.6. Theorem ([4, Theorem 6.4]) If $f: X \longrightarrow Y$ is a mapping of a compact $C$-space $X$ onto a space $Y$ such that $\left|f^{-1}(y)\right|<\mathfrak{c}$ for every $y \in Y$, then $Y$ is a $C$-space.

For each space $X$ and $n<\omega$ we let

$$
G_{n}(X)=\bigcup\{U \subset X: U \text { is open and } \operatorname{dim~ClU} \leq \mathrm{n}\}
$$

and

$$
S(X)=X-\bigcup_{n<\omega} G_{n}(X) .
$$

Sklyarenko ([9, Theorem 3]) proved the following lemma in the case when $X$ is S-weakly infinite dimensional.
2.7. Lemma $A$ weakly paracompact space $X$ is a finite $C$-space if and only if $S(X)$ is a compact finite $C$-space and every closed subspace $F \subset X$ disjoint from $S(X)$ is finite dimensional.

Proof. Assume that the space $X$ is a finite $C$-space. We shall show that $S(X)$ is compact. Suppose $S(X)$ is not compact. Since $S(X)$ is weakly paracompact, $S(X)$ is not pseudocompact. Thus there exists a countable discrete closed subspace $F$ of $S(X)$. Let us set $F=\left\{x_{i}: i<\omega\right\}$. We can take a discrete collection $\left\{U_{i}: i<\omega\right\}$ of open subsets of $X$ with $x_{i} \in U_{i}$ for every $i<\omega$. Thus we have $\operatorname{dim} \mathrm{Cl} U_{i}>i$ for every $i<\omega$. Let us set $Y=\bigcup\left\{\mathrm{Cl} U_{i}: i<\omega\right\}$. Since $\bigcup\left\{\mathrm{Cl} U_{i}: i<\omega\right\}$ is homeomorphic to $\bigoplus\left\{\mathrm{Cl} U_{i}: i<\omega\right\}, Y$ is not a S-weakly infinite dimensional subspace of $X$. Thus $Y$ is not a finite $C$-space. The contradiction shows that $S(X)$ is compact. Let F be a closed subset of $X$ disjoint from $S(X)$. First, we shall show that $F \subset G_{n}(X)$ for some $n<\omega$. Suppose that for every $n<\omega, F \not \subset G_{n}(X)$. Since $F \backslash G_{n}(X)$ is infinite for every $n<\omega$, inductively, we choose points $x_{1}, x_{2}, \cdots$ such that $x_{n} \in F \backslash\left(G_{n}(X) \cup\left\{x_{1}, x_{2}, \cdots, x_{n-1}\right\}\right)$ for every $n<\omega$. The space $E=\left\{x_{n}: n<\omega\right\}$ is a closed discrete subspace of $F$. For a space $E=\left\{x_{n}: n<\omega\right\}$ repeating above procedure we obtain a contradiction. Thus $F \subset G_{n}(X)$ for some $n<\omega$. Since $X$ is weakly paracompact, by the point finite sum theorem, $\operatorname{dim} F \leq n$. By Lemma B, the converse holds. Lemma 2.7 has been proved.
2.8 Proof of Theorem 2.5. By Lemma 2.7, $S(X)$ is compact. Applying Theorem 2.6 to $\left.f\right|_{S(X)}, f(S(X))$ is a finite $C$-space. For each closed subspace $F \subset Y$ disjoint from $f\left(S(X)\right.$ ), as $f^{-1}(F) \cap S(X)=\emptyset$, by Lemma 2.7, we take an integer $n$ with $\operatorname{dim} f^{-1}(F) \leq n$. As the restriction $\left.f\right|_{f^{-1}(F)}: f^{-1}(F) \longrightarrow F$ is open, by Nagami [7] (cf.
[3, 3.3.G]), $\operatorname{dim} F=\operatorname{dim} f^{-1}(F) \leq n$. Thus $F$ is a finite $C$-space. By Lemma B, $Y$ is a finite $C$-space.

## 3 Dimension-lowering mappings

The following theorem is a counterpart for $C$-spaces of Polkowski's result, which was proved in the case when A-weakly infinite-dimensional (see [8, Theorem 3.3 (ii)]).
3.1. Theorem If $f: X \longrightarrow Y$ is an open mapping of a paracompact space $X$ onto a $C$-space $Y$ such that $\left|f^{-1}(y)\right|<\omega$ for every $y \in Y$, then $X$ is a $C$-space.

To prove Theorem 3.1 we need the following lemma.
3.2. Lemma( $[7]$, cf $[8$, Lemma B]) If $f: X \longrightarrow Y$ is an open mapping of a space $X$ to a space $Y$ and there exists an integer $n \geq 1$ such that $\left|f^{-1}(y)\right|=n$ for every $y \in Y$, then $f$ is a local homeomorphism.
3.3 Proof of Theorem 3.1. For every $n \in \mathbb{N}$ we set

$$
Y_{n}=\left\{y \in Y:\left|f^{-1}(y)\right|=n\right\} \text { and } X_{n}=f^{-1}\left(Y_{n}\right) .
$$

It is easy to see that the union $Y_{n}^{\prime}=\bigcup_{k \leq n} Y_{k}$ is closed in $Y$ for every $n \in \mathbb{N}$, therefore the union $X_{n}^{\prime}=\bigcup_{k \leq n} X_{k}$ is also closed in $X$. Since $X$ is the union of countable collection $\left\{X_{n}^{\prime}: n \in \mathbb{N}\right\}$ of closed subsets of $X$, by the countable sum theorem for $C$-spaces, we only prove that $X_{n}^{\prime}$ is a $C$-space for every $n \in \mathbb{N}$. Let $f_{n}: X_{n} \longrightarrow Y_{n}$ be the mapping defined by $f_{n}(x)=f(x)$ for every $x \in X_{n}$.

Obviously, $X_{1}^{\prime}$ is a $C$-space, because $f_{1}$ is a homeomorphism. Assume that $X_{n-1}^{\prime}$ is a $C$-space. To prove that $X_{n}^{\prime}$ is a $C$-space, it suffices to show that every closed subset $Z$ of $X_{n}^{\prime}$ contained in $X_{n}$ is a $C$-space.

By Lemma 3.2, the mapping $f_{n}$ is a local homeomorphism. Thus for every $x \in X_{n}$ we can take a neighborhood $U_{x}$ of $x$ in $X_{n}$ such that the restriction $\left.f_{n}\right|_{U_{x}}: U_{x} \longrightarrow Y_{n}$ is an embedding. Since $X_{n}$ is open in $X_{n}^{\prime}, U_{x}$ is open in $X_{n}^{\prime}$. We may assume that $U_{x}$ is an $F_{\sigma}$-set of $X_{n}^{\prime}$. Let $U_{x}=\cup\{A(x, m): m \in \mathbb{N}\}$, where $A(x, m)$ is closed in $X_{n}^{\prime}$. For
every $y \in Y_{n}$ let us set $f^{-1}(y)=\{x(y, 1), x(y, 2), \ldots, x(y, n)\}$. Then the intersection $\bigcap_{i=1}^{n} f\left(U_{x(y, i)}\right)$ is a neighborhood of $y$ in $Y_{n}^{\prime}$. Take an open $F_{\sigma}$-set $V_{y}$ of $y$ in $Y_{n}^{\prime}$ such that $y \in V_{y} \subset \bigcap_{i=1}^{n} f\left(U_{x(y, i)}\right)$. Let $V_{y}=\cup\{B(y, \ell): \ell \in \mathbb{N}\}$, where $B(y, \ell)$ is closed in $Y_{n}^{\prime}$. The set $W(y, i)=U_{x(y, i)} \cap f^{-1}\left(V_{y}\right)$ is homeomorphic to $f(W(y, i))$. We have

$$
W(y, i)=\bigcup\left\{A(x(y, i), m) \cap f^{-1}(B(y, \ell)): m, \ell \in \mathbb{N}\right\} .
$$

We shall prove that $A(x(y, i), m) \cap f^{-1}(B(y, \ell))$ is a $C$-space. Since $\left.f_{n}\right|_{U_{x}(y, i)}$ is an embedding, $A(x(y, i), m) \cap f^{-1}(B(y, \ell))$ is homeomorphic to $f_{n}\left(A(x(y, i), m) \cap f^{-1}(B(y, \ell))\right)$.

By Lemma 2.3, $f_{n}$ is closed, therefore $f_{n}\left(A(x(y, i), m) \cap f^{-1}(B(y, \ell))\right)$ is closed in $Y_{n}$. Since $f_{n}\left(A(x(y, i), m) \cap f^{-1}(B(y, \ell))\right) \subset B(y, \ell) \subset Y_{n}, f_{n}\left(A(x(y, i), m) \cap f^{-1}(B(y, \ell))\right)$ is closed in $B(y, \ell)$. As $B(y, \ell)$ is a $C$-space, $f_{n}\left(A(x(y, i), m) \cap f^{-1}(b(y, \ell))\right)$ is a $C$ space. Thus $A(x(y, i), m) \cap f^{-1}(b(y, \ell))$ is a $C$-space. By the countable sum theorem for $C$-spaces, $W(y, i)$ is a $C$-space. Since $Z$ is paracompact, the open cover $\mathcal{W}=$ $\left\{W(y, i) \cap Z: y \in Y_{n}, 1 \leq i \leq n\right\}$ of $Z$ has a locally-finite closed refinement $\mathcal{F}$. Since every member of $\mathcal{F}$ is a $C$-space, by the locally finite sum theorem for $C$-spaces (cf. [6, Theorem 1.1(i)]), $Z$ is a $C$-space. Theorem 3.1 has been proved.
3.4. Theorem If $f: X \longrightarrow Y$ is a closed-and-open mapping of a space $X$ onto a weakly paracompact finite $C$-space $Y$ such that $\left|f^{-1}(y)\right|<\omega$ for every $y \in Y$, then $X$ is a finite $C$-space.

Proof. Since for every $y \in Y\left|f^{-1}(y)\right|<\omega$, the closed mapping $f: X \longrightarrow Y$ is perfect. As $S(Y)$ is compact, $f^{-1}(S(Y))$ is compact. By Theorem 3.1, $f^{-1}(S(Y))$ is a finite $C$-space. For each closed subset $F \subset X$ disjoint from $f^{-1}(S(Y))$, as $f(F) \cap S(Y)=\emptyset$, by Lemmma 2.7, we take integer $n$ with $\operatorname{dim} f(F) \leq n$. As $\left.f\right|_{F}: F \longrightarrow f(F)$ is closed, by [3, Theorem 3.3.10], $\operatorname{dim} F \leq \operatorname{dim} f(F) \leq n$. Thus $F$ is a finite $C$-space, by Lemma $\mathrm{B}, X$ is a finite $C$-space.

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