# Acta Universitatis Sapientiae 

Mathematica<br>Volume 5, Number 1, 2013

Dedicated to the memory of Professor Antal Bege (1962-2012), first Editor-in-Chief of Acta Universitatis Sapientiae, first Executive editor of series Mathematica

Sapientia Hungarian University of Transylvania Scientia Publishing House

## Contents

Sz. András
Data dependence of solutions for Fredholm-Volterra integral equations in $L^{2}[a, b]$ ..... 5
M. Bencze, M. Drăgan
Some inequalities in bicentric quadrilateral ..... 20
Z. Finta
Approximation by limit $q$-Bernstein operator ..... 39
E. Horobȩ
Galois covering and smash product of skew categories ..... 47
A. Iványi
Leader election in synchronous networks ..... 54
R. Szász, P. Kupán
About a condition for starlikeness ..... 83
L. Tóth
A survey of the alternating sum-of-divisors function ..... 93

# Data dependence of solutions for Fredholm-Volterra integral equations in $L^{2}[a, b]$ 

Szilárd András<br>Babeş-Bolyai University<br>Department of Applied Mathematics<br>Cluj-Napoca, M. Kogălniceanu, No. 1, Romania<br>email: andrasz@math.ubbcluj.ro

Dedicated to the memory of Professor Antal Bege


#### Abstract

In this paper we study the continuous dependence and the differentiability with respect to the parameter $\lambda \in\left[\lambda_{1}, \lambda_{2}\right]$ of the solution operator $S:\left[\lambda_{1}, \lambda_{2}\right] \rightarrow L^{2}[a, b]$ for a mixed Fredholm-Volterra type integral equation. The main tool is the fiber Picard operators theorem (see [9], [8], [11], [3] and [2]).


## 1 Introduction

We study the solution operator of the equations

$$
\begin{equation*}
y(x)=f(x)+\int_{a}^{x} K_{1}(x, s, y(s) ; \lambda) d s+\int_{a}^{b} K_{2}(x, s, y(s) ; \lambda) d s \tag{1}
\end{equation*}
$$

and

$$
\begin{equation*}
y(x)=f(x)+\int_{a}^{x} K_{1}\left(x, s, y\left(g_{1}(s)\right) ; \lambda\right) d s+\int_{a}^{b} K_{2}\left(x, s, y\left(g_{2}(s)\right) ; \lambda\right) d s \tag{2}
\end{equation*}
$$

2010 Mathematics Subject Classification: 47H10, 45G15
Key words and phrases: Fredholm-Volterra integral equations, fiber Picard operators
where $\lambda \in\left[\lambda_{1}, \lambda_{2}\right]$ is a real parameter. The existence and uniqueness of the solutions of such equations in $C[a, b]$ was studied by many authors [5], [6], [1], we recall the results from [1]. If the functions $K_{i}$ and $f$ satisfy the conditions under which the existence and uniqueness (in $C[a, b]$ ) is guaranteed then the differentiability of the functions $\mathrm{K}_{i}$ with respect to the parameter guarantees the differentiability of the solution. This property was proved in [1] using the following fiber Picard operator theorem:

Theorem 1 (Fiber Picard operator's) [9] Let (V, d) be a generalized metric space with $\mathrm{d}\left(\nu_{1}, \nu_{2}\right) \in \mathbb{R}_{+}^{p}$, and $(\mathbf{W}, \rho)$ a complete generalized metric space with $\rho\left(\mathcal{w}_{1}, w_{2}\right) \in \mathbb{R}_{+}^{\mathrm{m}}$. Let $\mathrm{A}: \mathrm{V} \times \mathrm{W} \rightarrow \mathrm{V} \times \mathrm{W}$ be a continuous operator. If we suppose that:
a) $\mathrm{A}(v, w)=(\mathrm{B}(v), \mathrm{C}(v, w))$ for all $v \in \mathrm{~V}$ and $w \in \mathrm{~W}$;
b) the operator $\mathrm{B}: \mathrm{V} \rightarrow \mathrm{V}$ is a weakly Picard operator;
c) there exists a matrix $\mathrm{Q} \in \mathrm{M}_{\mathrm{m}}\left(\mathbb{R}_{+}\right)$convergent to zero, such that the operator $\mathrm{C}(\nu, \cdot): \mathrm{W} \rightarrow \mathrm{W}$ is a Q contraction for all $\nu \in \mathrm{V}$,
then the operator A is a weakly Picard operator. Moreover, if B is a Picard operator, then the operator $\mathcal{A}$ is a Picard operator.

In this paper we use the same technique to give some modified Carathéodory type conditions which guarantee the continuity and differentiability with respect to the parameter of the solution operator. We study these equations both in bounded and unbounded intervals.

## 2 Fredholm-Volterra equations on a compact interval

We need the following lemma.
Lemma 1 If $\mathrm{I}=[\mathrm{a}, \mathrm{b}], \mathrm{k} \in \mathrm{L}^{2}\left(\mathrm{I}^{2}\right)$ and the function $u \in \mathrm{~L}^{2}(\mathrm{I})$ has nonnegative values then the inequality

$$
\begin{equation*}
u(t) \leq \alpha+\int_{a}^{b} k(t, s) u(s) d s, \text { a.e. } t \in I \tag{3}
\end{equation*}
$$

where $\alpha>0$ and $\|\mathrm{k}\|_{\mathrm{L}^{2}\left(\mathrm{I}^{2}\right)}<1$, implies

$$
\|u\|_{\mathrm{L}^{2}(\mathrm{I})} \leq \frac{\alpha \sqrt{2(\mathrm{~b}-\mathrm{a})}}{1-\|\mathrm{k}\|_{\mathrm{L}^{2}\left(\mathrm{I}^{2}\right)}}
$$

Proof. Consider the sets

$$
A=\{t \in I \mid u(t) \leq \alpha\} \text { and } B=\{t \in I \mid u(t)>\alpha\}
$$

These sets are measurable because $u$ is measurable. If $t \in B$, from the CauchyBuniakovski inequality we have

$$
(u(t)-\alpha)^{2} \leq\left(\int_{a}^{b} k(t, s) u(s) d s\right)^{2} \leq \int_{a}^{b} k^{2}(t, s) d s \cdot \int_{a}^{b} u^{2}(s) d s
$$

By integrating on $B$ we deduce

$$
\begin{aligned}
\int_{B} u^{2}(s) d s & \leq 2 \alpha \int_{B} u(t) d t-\alpha^{2} \cdot \mu(B)+\int_{B} \int_{a}^{b} k^{2}(t, s) d s d t \cdot\|u\|_{L^{2}(I)}^{2} \\
& \leq 2 \alpha \int_{B} u(t) d t-\alpha^{2} \cdot \mu(B)+\int_{a}^{b} \int_{a}^{b} k^{2}(t, s) d s d t \cdot\|u\|_{L^{2}(I)}^{2} \\
& \leq 2 \alpha \sqrt{\mu(B) \int_{a}^{b} u^{2}(t) d t}-\alpha^{2} \cdot \mu(B)+\|k\|_{L^{2}\left(I^{2}\right)}^{2} \cdot\|u\|_{L^{2}(I)}^{2}
\end{aligned}
$$

By the other hand $u^{2}(t) \leq \alpha^{2}$, for $t \in A$, so

$$
\int_{A} u^{2}(t) d t \leq \alpha^{2} \cdot \mu(A)
$$

From these inequalities we have

$$
\left(\|u\|_{\mathrm{L}^{2}(\mathrm{I})}-\alpha \sqrt{\mu(\mathrm{B})}\right)^{2} \leq \alpha^{2} \mu(\mathrm{~A})+\|\mathrm{k}\|_{\mathrm{L}^{2}\left(\mathrm{I}^{2}\right)}^{2} \cdot\|u\|_{\mathrm{L}^{2}(\mathrm{I})}^{2}
$$

so

$$
\|u\|_{\mathrm{L}^{2}(\mathrm{I})}-\alpha \sqrt{\mu(\mathrm{B})} \leq \sqrt{\alpha^{2} \mu(\mathrm{~A})+\|\mathrm{k}\|_{\mathrm{L}^{2}\left(\mathrm{I}^{2}\right)}^{2} \cdot\|u\|_{\mathrm{L}^{2}(\mathrm{I})}^{2}}
$$

From

$$
\sqrt{\alpha^{2} \mu(A)+\|\mathrm{k}\|_{\mathrm{L}^{2}\left(\mathrm{I}^{2}\right)}^{2} \cdot\|u\|_{\mathrm{L}^{2}(\mathrm{I})}^{2}} \leq \alpha \sqrt{\mu(A)}+\|\mathrm{k}\|_{\mathrm{L}^{2}\left(\mathrm{I}^{2}\right)} \cdot\|u\|_{\mathrm{L}^{2}(\mathrm{I})}
$$

and

$$
\sqrt{\mu(A)}+\sqrt{\mu(B)} \leq \sqrt{2(b-a)}
$$

we deduce the desired inequality.
Remark 1 By using both the Minkovski and the Cauchy-Buniakovski inequality we can prove a sharpened version:

$$
\|u\|_{L^{2}(\mathrm{I})} \leq \frac{\alpha \sqrt{(\mathrm{b}-\mathrm{a})}}{1-\|\mathrm{k}\|_{\mathrm{L}^{2}\left(\mathrm{I}^{2}\right)}}
$$

Indeed (3) implies
$\|u\|_{L^{2}(I)} \leq \| \alpha+\sqrt{\int_{a}^{b} k^{2}(t, s) d s \cdot \int_{a}^{b} u^{2}(s) d s\left\|_{L^{2}(I)} \leq \alpha \sqrt{b-a}+\right\| k\left\|_{L^{2}\left(I^{2}\right)} \cdot\right\| u \|_{L^{2}(I)} .}$
By an analogous reasoning we have the following property: If $\mathrm{k} \in \mathrm{L}^{2}\left(\mathrm{I}^{2}\right), \mathrm{g} \in$ $\mathrm{L}^{2}(\mathrm{I})$ and the function $u \in \mathrm{~L}^{2}(\mathrm{I})$ has nonnegative values then the inequality

$$
u(t) \leq g(t)+\int_{a}^{b} k(t, s) u(s) d s, \text { a.e. } t \in I
$$

where $\|\mathrm{k}\|_{\mathrm{L}^{2}\left(\mathrm{I}^{2}\right)}<1$, implies

$$
\|u\|_{\mathrm{L}^{2}(\mathrm{I})} \leq \frac{\|\mathrm{g}\|_{\mathrm{L}^{2}(\mathrm{I})}}{1-\|\mathrm{k}\|_{\mathrm{L}^{2}\left(\mathrm{I}^{2}\right)}}
$$

These inequalities are in fact Gronwall type inequalities and they can be proved also by using the abstract Gronwall lemma from [10].

We use the usual definition of differentiability for functions with values in a Banach space and a generalized Weierstrass type theorem. To avoid any misunderstanding we recall this definition and we prove the above mentioned theorem.

Definition 1 If $\mathrm{S}:\left[\lambda_{1}, \lambda_{2}\right] \rightarrow \mathrm{L}^{2}(\mathrm{I})$ is a continuous function then we call it differentiable at the point $\lambda$, if exists $z_{\lambda} \in \mathrm{L}^{2}(\mathrm{I})$ such that

$$
\lim _{\bar{\lambda} \rightarrow \lambda} \frac{\left\|S(\bar{\lambda})-S(\lambda)-(\bar{\lambda}-\lambda) z_{\lambda}\right\|_{L^{2}(\mathrm{I})}}{\bar{\lambda}-\lambda}=0
$$

For the simplicity we identify the function $\mathrm{t} \rightarrow \mathrm{t} \boldsymbol{z}_{\lambda}$ (the differential) with the element $z_{\lambda}$.

Theorem 2 If the sequence $\mathrm{y}_{\mathfrak{n}}(\cdot, \lambda) \in \mathrm{L}^{2}(\mathrm{I}), \mathrm{n} \geq 0$ converges in $\mathrm{L}^{2}(\mathrm{I})$ to $y^{*}(\cdot, \lambda)$ for all $\lambda \in\left[\lambda_{1}, \lambda_{2}\right]$, the operators $S_{n}:\left[\lambda_{1}, \lambda_{2}\right] \rightarrow L^{2}(I)$ defined by $\mathrm{S}_{\mathrm{n}}(\lambda)(\mathrm{t})=\mathrm{y}_{\mathrm{n}}(\mathrm{t}, \lambda), \forall \mathrm{t} \in \mathrm{I}, \forall \lambda \in\left[\lambda_{1}, \lambda_{2}\right]$ are differentiable, their differentials converge in $\mathrm{L}^{2}(\mathrm{I})$ to $z^{*}(\cdot, \lambda)$, and these convergencies are uniform with respect to $\lambda$, then the operator $S:\left[\lambda_{1}, \lambda_{2}\right] \rightarrow \mathrm{L}^{2}(\mathrm{I})$ defined by $\mathrm{S}(\lambda)(\mathrm{t})=\mathrm{y}^{*}(\mathrm{t}, \lambda), \forall \mathrm{t} \in$ I, $\forall \lambda \in\left[\lambda_{1}, \lambda_{2}\right]$ is differentiable and $z^{*}(\cdot, \lambda)$ is its differential in $\lambda$.

Proof. Due to the mean theorem for functions with values in a Banach space (see [4] 2-5) we have the following inequality:

$$
\begin{aligned}
& \left\|\left[y_{m}(\cdot, \bar{\lambda})-y_{\mathfrak{n}}(\cdot, \bar{\lambda})\right]-\left[y_{m}(\cdot, \lambda)-y_{\mathfrak{n}}(\cdot, \lambda)\right]\right\|_{L^{2}(\mathrm{I})} \\
& \bar{\lambda}-\lambda \\
& \leq \sup _{\lambda \in\left[\lambda_{1}, \lambda_{2}\right]}\left\|z_{\mathfrak{m}}(\cdot, \lambda)-z_{\mathfrak{n}}(\cdot, \lambda)\right\|_{L^{2}(\mathrm{I})}
\end{aligned}
$$

where $z_{\mathfrak{m}}(\cdot, \lambda)$ is the differential of $S_{n}(\lambda)(\cdot)$.
The condition $\left\|z_{\mathfrak{n}}(\cdot, \lambda)-z^{*}(\cdot, \lambda)\right\|_{L^{2}(\mathrm{I})} \rightarrow 0$ uniform with respect to $\lambda$, implies that for every $\varepsilon>0$ exists $n_{1}(\varepsilon) \in \mathbb{N}$ such that

$$
\begin{equation*}
\frac{\left\|\left[y^{*}(\cdot, \bar{\lambda})-y^{*}(\cdot, \lambda)\right]-\left[y_{n}(\cdot, \bar{\lambda})-y_{n}(\cdot, \lambda)\right]\right\|_{L^{2}(\mathrm{I})}}{\bar{\lambda}-\lambda} \leq \frac{\varepsilon}{3}, \forall \mathrm{n} \geq \mathfrak{n}_{1}(\varepsilon) \tag{4}
\end{equation*}
$$

By the other hand for all $\varepsilon>0$ exists $n_{2}(\varepsilon) \in \mathbb{N}$ such that

$$
\begin{equation*}
\left\|z_{\mathrm{n}}(\cdot, \lambda)-z^{*}(\cdot, \lambda)\right\|_{\mathrm{L}^{2}(\mathrm{I})} \leq \frac{\varepsilon}{3}, \forall \mathrm{n} \geq \mathrm{n}_{2}(\varepsilon) \tag{5}
\end{equation*}
$$

and there exists $\delta>0$ such that

$$
\begin{equation*}
\frac{\left\|y_{n}(\cdot, \bar{\lambda})-y_{n}(\cdot, \lambda)-(\bar{\lambda}-\lambda) z_{n}(\cdot, \lambda)\right\|_{L^{2}(\mathrm{I})}}{\bar{\lambda}-\lambda} \leq \frac{\varepsilon}{3} \tag{6}
\end{equation*}
$$

if $|\bar{\lambda}-\lambda|<\delta$. From these relations we deduce

$$
\lim _{\bar{\lambda} \rightarrow \lambda} \frac{\left\|y^{*}(\cdot, \bar{\lambda})-y^{*}(\cdot, \lambda)-(\bar{\lambda}-\lambda) z^{*}(\cdot, \lambda)\right\|_{L^{2}(\mathrm{I})}}{\bar{\lambda}-\lambda}=0
$$

so $S$ is differentiable in $\lambda$ and its differential is $z^{*}(\cdot, \lambda)$.
For equation (1) we have the following theorem (some parts of this theorem are classical):

Theorem 3 If
I. (Carathéodory type conditions) the functions $\mathrm{K}_{\mathrm{i}}: \mathrm{I}^{2} \times\left[\lambda_{1}, \lambda_{2}\right] \times \mathbb{R} \rightarrow \mathbb{R}$, $\mathfrak{i} \in\{1,2\}$ with $\mathrm{I}=[\mathrm{a}, \mathrm{b}]$ satisfy the conditions
a ) $\mathrm{K}_{\mathfrak{i}}(\cdot, \cdot, \lambda, \mathfrak{u})$ is measurable on $\mathrm{I}^{2}=[\mathrm{a}, \mathrm{b}] \times[\mathrm{a}, \mathrm{b}]$ for all $\mathrm{u} \in \mathbb{R}$ and $\lambda \in\left[\lambda_{1}, \lambda_{2}\right] ;$
b) $\mathrm{K}_{\mathfrak{i}}(\mathrm{x}, \mathrm{s}, \lambda, \cdot)$ is continuous on $\mathbb{R}$ for almost every pairs $(\mathrm{x}, \mathrm{s}) \in \mathrm{I}^{2}$ and every $\lambda \in\left[\lambda_{1}, \lambda_{2}\right]$.
II. (space invariance) $\mathrm{f} \in \mathrm{L}^{2}(\mathrm{I}), \mathrm{K}_{\mathfrak{i}}(\cdot, \cdot, \lambda, 0) \in \mathrm{L}^{2}\left(\mathrm{I}^{2}\right)$ for all $\lambda \in\left[\lambda_{1}, \lambda_{2}\right]$, $\mathfrak{i} \in\{1,2\}$ and exists $M_{1}>0$ such that $\left\|\mathrm{K}_{\mathfrak{i}}(\cdot, \cdot, \lambda, 0)\right\|_{\mathrm{L}^{2}\left(\mathrm{I}^{2}\right)}<\mathrm{M}_{1}$ for all $\lambda \in\left[\lambda_{1}, \lambda_{2}\right] ;$
III. (Lipschitz type conditions) exists $\mathrm{k}_{\mathfrak{i}} \in \mathrm{L}^{2}\left(\mathrm{I}^{2}\right), \mathfrak{i} \in\{1,2\}$, such that
$\left|\mathrm{K}_{\mathfrak{i}}(\mathrm{t}, \mathrm{s}, \lambda, \mathrm{u})-\mathrm{K}_{\mathfrak{i}}(\mathrm{t}, \mathrm{s}, \lambda, v)\right| \leq \mathrm{k}_{\mathfrak{i}}(\mathrm{t}, \mathrm{s})|\mathrm{u}-v|, \forall \mathrm{t}, \mathrm{s} \in \mathrm{I}, \lambda \in\left[\lambda_{1}, \lambda_{2}\right], u, v \in \mathbb{R} ;$
IV. (contraction condition)

$$
\begin{equation*}
L^{2}:=\int_{a}^{b} \int_{a}^{t}\left(k_{1}(t, s)+k_{2}(t, s)\right)^{2} d s d t+\int_{a}^{b} \int_{t}^{b} k_{2}^{2}(t, s) d s d t<1 \tag{7}
\end{equation*}
$$

then

1. for all $\lambda \in\left[\lambda_{1}, \lambda_{2}\right]$ exists a unique solution $\mathrm{y}^{*}(\cdot, \lambda) \in \mathrm{L}^{2}(\mathrm{I})$ of the equation (1);
2. the sequence of successive approximation

$$
y_{n+1}(x)=f(x)+\int_{a}^{x} K_{1}\left(x, s, \lambda, y_{n}(s)\right) d s+\int_{a}^{b} K_{2}\left(x, s, \lambda, y_{n}(s)\right) d s
$$

converges in $\mathrm{L}^{2}(\mathrm{I})$ to $\mathrm{y}^{*}(\cdot, \lambda)$, for all $\mathrm{y}_{0}(\cdot) \in \mathrm{L}^{2}(\mathrm{I})$ and every $\lambda \in\left[\lambda_{1}, \lambda_{2}\right]$;
3. for every $\mathfrak{n} \in \mathbb{N}$ we have

$$
\left\|y_{n}(\cdot)-y^{*}(\cdot, \lambda)\right\|_{L^{2}(\mathrm{I})} \leq \frac{\mathrm{L}^{n}}{1-\mathrm{L}}\left\|\mathrm{y}_{1}(\cdot)-\mathrm{y}_{0}(\cdot)\right\|_{\mathrm{L}^{2}(\mathrm{I})}
$$

## Moreover if

I.c) the functions $\left(\mathrm{K}_{\mathfrak{i}}(\mathrm{x}, \mathrm{s}, \cdot, \mathrm{u})\right)_{\mathrm{x}, \mathrm{s} \in \mathrm{I}, \mathrm{u} \in \mathbb{R}}$ are equally continuous,
then the operator $S:\left[\lambda_{1}, \lambda_{2}\right] \rightarrow L^{2}(I)$ defined by $S(\lambda)(x)=y^{*}(x, \lambda), \forall x \in I$, $\forall \lambda \in\left[\lambda_{1}, \lambda_{2}\right]$ is continuous.

If instead of I.b), I.c) and III. we have the conditions
I.b') $\mathrm{K}_{\mathfrak{i}}(\mathrm{x}, \mathrm{s}, \lambda, \cdot)$ is in $\mathrm{C}^{1}(\mathbb{R})$ for all $\lambda \in\left[\lambda_{1}, \lambda_{2}\right]$, a.e. $(\mathrm{x}, \mathrm{s}) \in \mathrm{I}^{2}$, and there exist $\mathrm{k}_{\mathrm{i}} \in \mathrm{L}^{2}\left(\mathrm{I}^{2}\right), \mathfrak{i} \in\{1,2\}$, such that

$$
\left|\frac{\partial \mathrm{K}_{\mathrm{i}}(\mathrm{t}, \mathrm{~s}, \lambda, \mathrm{u})}{\partial \mathrm{u}}\right| \leq \mathrm{k}_{\mathrm{i}}(\mathrm{t}, \mathrm{~s}), \forall \mathrm{t}, \mathrm{~s} \in \mathrm{I}, \forall \lambda \in\left[\lambda_{1}, \lambda_{2}\right], \forall u \in \mathbb{R}
$$

I.c') $\mathrm{K}_{\mathfrak{i}}(\mathrm{x}, \mathrm{s}, \cdot, \mathrm{u})$ is in $\mathrm{C}^{1}\left[\lambda_{1}, \lambda_{2}\right]$ for all $\mathrm{u} \in \mathbb{R}$, a.e. $(\mathrm{x}, \mathrm{s}) \in \mathrm{I}^{2}$, the partial derivatives satisfy condition $\mathrm{I} ., \frac{\partial \mathrm{K}_{\mathrm{i}}}{\partial \lambda}(\cdot, \cdot, \lambda, u) \in \mathrm{L}^{2}\left(\mathrm{I}^{2}\right), \mathfrak{i} \in\{1,2\}$ and there exists $\mathrm{M}_{2}>0$ such that

$$
\left\|\frac{\partial \mathrm{K}_{\mathfrak{i}}}{\partial \lambda}(\cdot, \cdot, \lambda, u)\right\|_{\mathrm{L}^{2}\left(\mathrm{I}^{2}\right)}<M_{2}, \forall \lambda \in\left[\lambda_{1}, \lambda_{2}\right], \forall u \in \mathbb{R}
$$

then the operator S is differentiable.
Proof. First we prove that for a fixed $\lambda$ the operator $T: L^{2}(I) \rightarrow L^{2}(I)$ defined by

$$
T[y](x)=f(x)+\int_{a}^{x} K_{1}(x, s, \lambda, y(s)) d s+\int_{a}^{b} K_{2}(x, s, \lambda, y(s)) d s
$$

is a contraction. From the Lipschitz condition we have

$$
\int_{a}^{b} K_{2}(t, s, \lambda, y(s)) d s \leq \int_{a}^{b} K_{2}(t, s, \lambda, 0)+k_{2}(t, s)|y(s)| d s
$$

Due to Minkovski and Cauchy-Buniakovski inequality we deduce

$$
\begin{gathered}
\int_{a}^{b}\left(\int_{a}^{b} K_{2}(t, s, \lambda, y(s)) d s\right)^{2} d t \\
\leq\left(\sqrt{b-a}\left\|K_{2}(\cdot, \cdot, \lambda, 0)\right\|_{L^{2}\left(I^{2}\right)}+\sqrt{b-a}\left\|k_{2}\right\|_{L^{2}\left(I^{2}\right)} \cdot\|y\|_{L^{2}(I)}\right)^{2}<\infty
\end{gathered}
$$

Analogously

$$
\int_{a}^{b}\left(\int_{a}^{t} K_{1}(t, s, \lambda, y(s)) d s\right)^{2} d t<\infty
$$

so because of $f \in L^{2}(I)$ we have $T[y] \in L^{2}(I)$. On the other hand

$$
\begin{aligned}
\left|T\left[y_{1}\right](t)-T\left[y_{2}\right](t)\right| \leq & \int_{a}^{t}\left|K_{1}\left(t, s, \lambda, y_{1}(s)\right)-K_{1}\left(t, s, \lambda, y_{2}(s)\right)\right| d s \\
& +\int_{a}^{b}\left|K_{2}\left(t, s, \lambda, y_{1}(s)\right)-K_{2}\left(t, s, \lambda, y_{2}(s)\right)\right| d s \\
\leq & \int_{a}^{t} k_{1}(t, s)\left|y_{1}(s)-y_{2}(s)\right| d s+\int_{a}^{b} k_{2}(t, s)\left|y_{1}(s)-y_{2}(s)\right| d s \\
= & \int_{a}^{b}\left(\bar{k}_{1}(t, s)+k_{2}(t, s)\right)\left|y_{1}(s)-y_{2}(s)\right| d s
\end{aligned}
$$

where $\bar{k}_{1}(t, s)=\left\{\begin{array}{ll}k_{1}(t, s), & t \geq s \\ 0, & t<s\end{array}\right.$. From the Cauchy-Buniakovski inequality we obtain

$$
\left\|\mathrm{T}\left[\mathrm{y}_{1}\right](\cdot)-\mathrm{T}\left[\mathrm{y}_{2}\right](\cdot)\right\|_{\mathrm{L}^{2}(\mathrm{I})}^{2} \leq \mathrm{L}^{2} \cdot\left\|\mathrm{y}_{1}(\cdot)-\mathrm{y}_{2}(\cdot)\right\|_{\mathrm{L}^{2}(\mathrm{I})}^{2},
$$

where $\mathrm{L}^{2}$ is defined by (7). Hence T is a contraction and from the contractions principle we have the conclusions.

If we have condition I.c), then for every $\varepsilon>0$ there exists $\varepsilon_{1}=\frac{(1-L) \varepsilon}{2(b-a) \sqrt{2(b-a)}}$ and $\delta>0$ such that for $|\lambda-\bar{\lambda}|<\delta$ we have

$$
\left|K_{i}(t, s, \lambda, u)-K_{i}(t, s, \bar{\lambda}, u)\right| \leq \varepsilon_{1},
$$

for all $u \in \mathbb{R}$ and a.e. $(t, s) \in I^{2}$. If $y_{\lambda}^{*}$ and $y_{\bar{\lambda}}^{*}$ are the corresponding unique solutions to $\lambda$, and $\bar{\lambda}$, then

$$
\begin{aligned}
\left|y_{\lambda}^{*}(t)-y_{\bar{\lambda}}^{*}(t)\right| \leq & \int_{a}^{t}\left|K_{1}\left(t, s, \lambda, y_{\lambda}^{*}(s)\right)-K_{1}\left(t, s, \bar{\lambda}, y_{\bar{\lambda}}^{*}(s)\right)\right| d s \\
& +\int_{a}^{b}\left|K_{2}\left(t, s, \lambda, y_{\lambda}^{*}(s)\right)-K_{2}\left(t, s, \bar{\lambda}, y_{\bar{\lambda}}^{*}(s)\right)\right| d s \\
\leq & 2(b-a) \varepsilon_{1}+\int_{a}^{t}\left|K_{1}\left(t, s, \lambda, y_{\lambda}^{*}(s)\right)-K_{1}\left(t, s, \lambda, y_{\bar{\lambda}}^{*}(s)\right)\right| d s \\
& +\int_{a}^{b}\left|K_{2}\left(t, s, \lambda, y_{\lambda}^{*}(s)\right)-K_{2}\left(t, s, \lambda, y_{\bar{\lambda}}^{*}(s)\right)\right| d s \\
\leq & 2(b-a) \varepsilon_{1}+\int_{a}^{b}\left(\bar{k}_{1}(t, s)+k_{2}(t, s)\right)\left|y_{\lambda}^{*}(s)-y_{\bar{\lambda}}^{*}(s)\right| d s .
\end{aligned}
$$

From this inequality and Lemma 1 we obtain

$$
\left\|y_{\lambda}^{*}(\cdot)-y_{\bar{\lambda}}^{*}(\cdot)\right\|_{L^{2}(\mathrm{I})} \leq \frac{2(b-a) \varepsilon_{1} \sqrt{2(b-a)}}{1-L}
$$

where $L$ is defined in (7). So for every $\varepsilon>0$ there exists $\delta>0$ such that

$$
|\lambda-\bar{\lambda}|<\delta \Rightarrow\left\|y_{\lambda}^{*}(\cdot)-y_{\lambda}^{*}(\cdot)\right\|_{\mathrm{L}^{2}(\mathrm{I})}<\varepsilon
$$

this is the continuity of the operator $S$.
If we have I.b') and I.c'), we use the fiber Picard theorem to study the differentiability of the operator $S$. Consider the spaces $V=W=L^{2}(I)$ and the
operators $\mathrm{B}: \mathrm{V} \rightarrow \mathrm{V}, \mathrm{C}: \mathrm{V} \times \mathrm{W} \rightarrow \mathrm{W}$ defined by the following relations

$$
B[v](t)=f(t)+\int_{a}^{t} K_{1}(t, s, \lambda, y(s)) d s+\int_{a}^{b} K_{2}(t, s, \lambda, y(s)) d s
$$

and

$$
\begin{aligned}
& C[(v, w)](t)=\int_{a}^{t} \frac{\partial K_{1}(t, s, v(s) ; \lambda)}{\partial \lambda} d s+\int_{a}^{b} \frac{\partial K_{2}(t, s, v(s) ; \lambda)}{\partial \lambda} d s \\
& +\int_{a}^{t} \frac{\partial K_{1}(t, s, v(s) ; \lambda)}{\partial v} w(s) d s+\int_{a}^{b} \frac{\partial K_{2}(t, s, v(s) ; \lambda)}{\partial v} w(s) d s .
\end{aligned}
$$

Due to the given conditions the operator B is a Picard operator (condition I.b') implies condition III.) and the operator C satisfies

$$
\left\|\mathrm{C}\left[\left(v, w_{1}\right)\right]-\mathrm{C}\left[\left(v, w_{2}\right)\right]\right\|_{\mathrm{L}^{2}(\mathrm{I})} \leq \mathrm{L}_{1}\left\|w_{1}-w_{2}\right\|_{\mathrm{L}^{2}(\mathrm{I})}
$$

where $L_{1}=\sqrt{\int_{a}^{b} a_{a}^{t}\left(k_{1}(t, s)+k_{2}(t, s)\right)^{2} d s d t+\int_{a}^{b} \int_{t}^{b} k_{2}^{2}(t, s) d s d t \text {. Theorem } 1110 .}$ implies that the triangular operator $A[v, w]=(B[v], C[v, w])$ is a Picard operator and so the sequence of successive approximations constructed by the relations $\left(y_{n+1}, z_{n+1}\right)=A\left[y_{n}, z_{n}\right]$ converges in $\left(L^{2}(I)\right)^{2}$ to the unique fixed point. If we choose for $y_{0}(\cdot, \lambda)$ a $C^{1}$ function in its last variable and $z_{0}=\frac{\partial y_{0}}{\partial \lambda}$, then from the definition of the operator $C$ we deduce $z_{n}=\frac{\partial y_{n}}{\partial \lambda}$. By the other hand the operators $S_{n}:\left[\lambda_{1}, \lambda_{2}\right] \rightarrow L^{2}(I)$ defined by $S_{n}(\lambda)(t)=y_{n}(t), \forall t \in$ I, $\forall \lambda \in\left[\lambda_{1}, \lambda_{2}\right]$ are differentiables and the differential of $S_{n}$ in $\lambda$ is $z_{n}$, hence we can apply Theorem 2 and we obtain the differentiability of the operator $S$.

Remark 2 We can prove the same results working in the space

$$
\mathrm{Y}=\left\{\mathrm{y}: \mathrm{I} \times \Lambda \rightarrow \mathbb{R} \mid \mathrm{y}(\cdot, \lambda) \in \mathrm{L}^{2}[\mathrm{I}], \forall \lambda \in \Lambda, \mathrm{y}(\mathrm{t}, \cdot) \in \mathrm{C}(\Lambda) \text { a.e. } \mathrm{t} \in \mathrm{I}\right\}
$$

where $\Lambda=\left[\lambda_{1}, \lambda_{2}\right]$ and the norm is defined by $\|y\|_{\mathrm{Y}}=\max _{\lambda \in \Lambda}\|\mathrm{y}(\cdot, \lambda)\|_{\mathrm{L}^{2}(\mathrm{I})}$.
Using the same arguments we can prove the following theorem for equation (2).

## Theorem 4 If

a) the functions $\mathrm{K}_{\mathrm{i}}: \mathrm{I} \times \mathrm{I} \times\left[\lambda_{1}, \lambda_{2}\right] \times \mathbb{R} \rightarrow \mathbb{R}, \mathfrak{i}=\overline{1,2}$ satisfy conditions I.-IV. from Theorem 3;
b) the functions $\mathrm{g}_{1}, \mathrm{~g}_{2}:[\mathrm{a}, \mathrm{b}] \rightarrow \mathbb{R}$ are injective and measurable and they satisfy the conditions $\operatorname{Im}\left(\mathrm{g}_{1}\right)=\left[\mathrm{a}_{1}, \mathrm{a}_{2}\right], \operatorname{Im}\left(\mathrm{g}_{2}\right)=\left[\mathrm{b}_{2}, \mathrm{~b}_{1}\right]$ with $\mathrm{a}_{1} \leq$ $\mathrm{a} \leq \mathrm{a}_{2} \leq \mathrm{b}$, and $\mathrm{a} \leq \mathrm{b}_{2} \leq \mathrm{b} \leq \mathrm{b}_{1}$;
c) $\varphi_{1} \in \mathrm{~L}^{2}\left(\left[a_{1}, a\right]\right)$ and $\varphi_{2} \in \mathrm{~L}^{2}\left(\left[\mathrm{~b}, \mathrm{~b}_{1}\right]\right)$;
then

1) equation (2) has a unique solution $\mathrm{y}^{*}(\cdot, \lambda)$ in $\mathrm{L}^{2}\left(\mathrm{I}_{1}\right)$ for all $\lambda \in\left[\lambda_{1}, \lambda_{2}\right]$, where $\mathrm{I}_{1}=\left[\mathrm{a}_{1}, \mathrm{~b}_{1}\right]$;
2) the sequence of successive approximations converges $L^{2}\left(I_{1}\right)$ to $y^{*}(\cdot, \lambda)$ for every admissible initial function $y_{0}(\cdot, \lambda)$, where the set of admissible functions is

$$
\begin{aligned}
\mathrm{Y}_{\mathrm{a}} & =\left\{\mathrm{y}(\cdot, \lambda) \in \mathrm{L}^{2}\left(\mathrm{I}_{1}\right) \mid y_{0}(\mathrm{t}, \lambda)=\varphi_{1}(\mathrm{t}), \forall \mathrm{t} \in\left[\mathrm{a}_{1}, \mathrm{a}\right], y_{0}(\mathrm{t}, \lambda)\right. \\
& \left.=\varphi_{2}(\mathrm{t}), \forall \mathrm{t} \in\left[\mathrm{~b}, \mathrm{~b}_{1}\right]\right\} ;
\end{aligned}
$$

3) we have the following estimation:

$$
\left\|y_{n}(\cdot)-y^{*}(\cdot, \lambda)\right\|_{L^{2}\left(I_{1}\right)} \leq \frac{L^{n}}{1-L}\left\|y_{1}(\cdot)-y_{0}(\cdot)\right\|_{L^{2}\left(I_{1}\right)}
$$

where L is defined by relation (7).
Moreover if condition I.c) holds, then the operator $S:\left[\lambda_{1}, \lambda_{2}\right] \rightarrow L^{2}\left(\mathrm{I}_{1}\right)$ defined by $S(\lambda)(x)=y^{*}(x, \lambda), \forall x \in\left[a_{1}, b_{1}\right], \forall \lambda \in\left[\lambda_{1}, \lambda_{2}\right]$ is continuous.

If instead of conditions I.b), I.c) and III. the conditions I.b') and I.c') are satisfied, then S is differentiable.

Remark 3 The differentiability of S implies the existence of the partial derivative $\frac{\partial y^{*}(\cdot, \lambda)}{\partial \lambda}$ and so from the construction of the operator C we deduce that this partial derivative satisfies the equation

$$
\begin{aligned}
& \frac{\partial y^{*}(t, \lambda)}{\partial \lambda}=\int_{a}^{t} \frac{\partial K_{1}\left(t, s, \lambda, y^{*}(s, \lambda)\right)}{\partial \lambda} d s+\int_{a}^{b} \frac{\partial K_{2}\left(t, s, \lambda, y^{*}(s, \lambda)\right)}{\partial \lambda} d s \\
& \quad+\int_{a}^{t} \frac{\partial K_{1}\left(t, s, \lambda, y^{*}(s, \lambda)\right)}{\partial y^{*}} \frac{\partial y^{*}(s, \lambda)}{\partial \lambda} d s+\int_{a}^{b} \frac{\partial K_{2}\left(t, s, \lambda, y^{*}(s, \lambda)\right)}{\partial y^{*}} \frac{\partial y^{*}(s, \lambda)}{\partial \lambda} d s
\end{aligned}
$$

in the case of Theorem 3 and the equation

$$
\begin{aligned}
\frac{\partial y^{*}(t, \lambda)}{\partial \lambda}= & \int_{a}^{t} \frac{\partial K_{1}\left(t, s, \lambda, y^{*}\left(g_{1}(s), \lambda\right)\right)}{\partial \lambda} d s+\int_{a}^{b} \frac{\partial K_{2}\left(t, s, \lambda, y^{*}\left(g_{2}(s), \lambda\right)\right)}{\partial \lambda} d s \\
& +\int_{a}^{t} \frac{\partial K_{1}\left(t, s, \lambda, y^{*}\left(g_{1}(s), \lambda\right)\right)}{\partial y^{*}} \cdot \frac{\partial y^{*}\left(g_{1}(s), \lambda\right)}{\partial \lambda} d s \\
& +\int_{a}^{b} \frac{\partial K_{2}\left(t, s, \lambda, y^{*}\left(g_{2}(s), \lambda\right)\right)}{\partial y^{*}} \cdot \frac{\partial y^{*}\left(g_{2}(s), \lambda\right)}{\partial \lambda} d s
\end{aligned}
$$

in the case of Theorem 4.

## 3 Fredholm-Volterra equations on an unbounded interval

If $I=[a, \infty)$, we can't use the same inequalities because in Lemma 1 and in some estimations we used it was essential the finite length of the interval. Due to this problem we need other conditions to guarantee the same properties of the solution operator.

Theorem 5 If conditions I.-III. from Theorem 3 are satisfied with $\mathrm{I}=[\mathrm{a}, \infty)$ and

$$
\begin{equation*}
L^{2}:=\int_{a}^{\infty} \int_{a}^{t}\left(k_{1}(t, s)+k_{2}(t, s)\right)^{2} d s d t+\int_{a}^{\infty} \int_{t}^{\infty} k_{2}^{2}(t, s) d s d t<1 \tag{8}
\end{equation*}
$$

then

1. for every $\lambda \in\left[\lambda_{1}, \lambda_{2}\right]$ there exists an unique solution $y^{*}(\cdot, \lambda) \in L^{2}(I)$;
2. the sequence of successive approximations

$$
y_{n+1}(x)=f(x)+\int_{a}^{x} K_{1}\left(x, s, \lambda, y_{n}(s)\right) d s+\int_{a}^{\infty} K_{2}\left(x, s, \lambda, y_{n}(s)\right) d s
$$

converges in $\mathrm{L}^{2}(\mathrm{I})$ to $\mathrm{y}^{*}(\cdot, \lambda)$, for all $\mathrm{y}_{0}(\cdot) \in \mathrm{L}^{2}(\mathrm{I})$;
3. for every $\mathfrak{n} \in \mathbb{N}$ we have

$$
\left\|y_{n}(\cdot)-y^{*}(\cdot, \lambda)\right\|_{L^{2}(\mathrm{I})} \leq \frac{\mathrm{L}^{n}}{1-\mathrm{L}}\left\|\mathrm{y}_{1}(\cdot)-y_{0}(\cdot)\right\|_{\mathrm{L}^{2}(\mathrm{I})}
$$

## Moreover if

I.c) there exist $\Lambda_{i}:\left[\lambda_{1}, \lambda_{2}\right] \times\left[\lambda_{1}, \lambda_{2}\right] \rightarrow \mathbb{R}$, and $g_{i}: I^{2} \rightarrow \mathbb{R}, \mathfrak{i} \in\{1,2\}$ such that
i)

$$
\begin{align*}
&\left|\mathrm{K}_{i}(\mathrm{x}, \mathrm{~s}, \lambda, u)-\mathrm{K}_{i}(\mathrm{x}, \mathrm{~s}, \bar{\lambda}, \mathrm{u})\right| \leq \Lambda_{i}(\lambda, \bar{\lambda}) \cdot \mathrm{g}_{i}(\mathrm{t}, \mathrm{~s}),  \tag{9}\\
& \forall u \in \mathbb{R}, \lambda, \bar{\lambda} \in\left[\lambda_{1}, \lambda_{2}\right], \text { a.e. }(\mathrm{t}, \mathrm{~s}) \in \mathrm{I}^{2}, \mathfrak{i} \in\{1,2\}
\end{align*}
$$

ii) $\lim _{\bar{\lambda} \rightarrow \lambda} \Lambda(\lambda, \bar{\lambda})=0$;
iii) $\int_{a}^{\infty}\left[\left(\int_{a}^{t} g_{1}(s, t) d s\right)^{2}+\left(\int_{a}^{\infty} g_{2}(s, t)\right)^{2}\right] d t<+\infty$
then the operator $S:\left[\lambda_{1}, \lambda_{2}\right] \rightarrow L^{2}(\mathrm{I})$ defined by $\mathrm{S}(\lambda)(x)=\mathrm{y}^{*}(\mathrm{x}, \lambda), \forall \mathrm{x} \in$ I, $\forall \lambda \in\left[\lambda_{1}, \lambda_{2}\right]$ is continuous.

If instead of the conditions I.b) and III. condition I.b') from Theorem 3 is fulfilled and
I.c') $\mathrm{K}_{\mathfrak{i}}(\mathrm{x}, \mathrm{s}, \cdot, \mathrm{u})$ is a $\mathrm{C}^{1}\left[\lambda_{1}, \lambda_{2}\right]$ function for all $u \in \mathbb{R}$, a.e. $(x, s) \in \mathrm{I}^{2}$, the partial derivatives satisfy condition I., and there exists $M_{3}>0$ such that

$$
\int_{a}^{\infty}\left(\int_{a}^{t} \frac{\partial K_{1}}{\partial \lambda}(t, s, \lambda, u) d s\right)^{2} d t+\int_{a}^{\infty}\left(\int_{a}^{t} \frac{\partial K_{2}}{\partial \lambda}(t, s, \lambda, u) d s\right)^{2} d t<M_{3}^{2}
$$

for all $\lambda \in\left[\lambda_{1}, \lambda_{2}\right]$ and for all $u \in \mathbb{R}$,
then S is differentiable.
Proof. As in Theorem 3 for a fixed $\lambda$ the operator $T: L^{2}(I) \rightarrow L^{2}(I)$ defined by

$$
T[y](x)=f(x)+\int_{a}^{x} K_{1}(x, s, \lambda, y(s)) d s+\int_{a}^{\infty} K_{2}(x, s, \lambda, y(s)) d s
$$

is a contraction with Lipschitz constant L. If $y_{\lambda}^{*}$ and $y_{\bar{\lambda}}^{*}$ are the unique solutions corresponding to $\lambda$ and $\bar{\lambda}$, from I.c) we deduce:

$$
\int_{a}^{\infty}\left(\int_{a}^{t}\left|K_{1}\left(t, s, \lambda, y_{\bar{\lambda}}^{*}(s)\right)-K_{1}\left(t, s, \bar{\lambda}, y_{\bar{\lambda}}^{*}(s)\right)\right| d s\right)^{2} d t \leq \Lambda_{1}^{2}(\lambda, \bar{\lambda}) \cdot \int_{a}^{\infty}\left(\int_{a}^{t} g_{1}(t, s) d s\right)^{2} d t
$$

and
$\int_{a}^{\infty}\left(\int_{a}^{\infty}\left|K_{2}\left(t, s, \lambda, y_{\bar{\lambda}}^{*}(s)\right)-K_{2}\left(t, s, \bar{\lambda}, y_{\bar{\lambda}}^{*}(s)\right)\right| d s\right)^{2} d t \leq \Lambda_{2}^{2}(\lambda, \bar{\lambda}) \cdot \int_{a}^{\infty}\left(\int_{a}^{\infty} g_{2}(t, s) d s\right)^{2} d t$.
From

$$
\begin{aligned}
\left|y_{\lambda}^{*}(t)-y_{\bar{\lambda}}^{*}(t)\right| \leq & \int_{a}^{t}\left|K_{1}\left(t, s, \lambda, y_{\lambda}^{*}(s)\right)-K_{1}\left(t, s, \bar{\lambda}, y_{\frac{*}{\lambda}}^{*}(s)\right)\right| d s \\
& +\int_{a}^{b}\left|K_{2}\left(t, s, \lambda, y_{\lambda}^{*}(s)\right)-K_{2}\left(t, s, \bar{\lambda}, y_{\frac{*}{\lambda}}^{*}(s)\right)\right| d s \\
\leq & \int_{a}^{t}\left|K_{1}\left(t, s, \lambda, y_{\bar{\lambda}}^{*}(s)\right)-K_{1}\left(t, s, \bar{\lambda}, y_{\bar{\lambda}}^{*}(s)\right)\right| d s \\
& +\int_{a}^{b}\left|K_{2}\left(t, s, \lambda, y_{\bar{\lambda}}^{*}(s)\right)-K_{2}\left(t, s, \bar{\lambda}, y_{\frac{*}{\lambda}}^{*}(s)\right)\right| d s \\
& +\int_{a}^{t}\left|K_{1}\left(t, s, \lambda, y_{\lambda}^{*}(s)\right)-K_{1}\left(t, s, \lambda, y_{\bar{\lambda}}^{*}(s)\right)\right| d s \\
& +\int_{a}^{b}\left|K_{2}\left(t, s, \lambda, y_{\lambda}^{*}(s)\right)-K_{2}\left(t, s, \lambda, y_{\bar{\lambda}}^{*}(s)\right)\right| d s \\
\leq & \int_{a}^{t}\left|K_{1}\left(t, s, \lambda, y_{\bar{\lambda}}^{*}(s)\right)-K_{1}\left(t, s, \bar{\lambda}, y_{\bar{\lambda}}^{*}(s)\right)\right| d s \\
& +\int_{a}^{b}\left|K_{2}\left(t, s, \lambda, y_{\bar{\lambda}}^{*}(s)\right)-K_{2}\left(t, s, \bar{\lambda}, y_{\bar{\lambda}}^{*}(s)\right)\right| d s \\
& +\int_{a}^{b}\left(\bar{k}_{1}(t, s)+k_{2}(t, s)\right)\left|y_{\lambda}^{*}(s)-y_{\bar{\lambda}}^{*}(s)\right| d s
\end{aligned}
$$

we deduce (using Minkovski inequality)

$$
\left\|y_{\lambda}^{*}(\cdot)-y_{\lambda}^{*}(\cdot)\right\|_{\mathrm{L}^{2}(\mathrm{I})} \leq \frac{\Lambda}{1-\mathrm{L}}
$$

where $L$ is defined in (8) and

$$
\Lambda=\Lambda_{1}(\lambda, \bar{\lambda}) \sqrt{\int_{a}^{\infty}\left(\int_{a}^{t} k_{1}(s, t) d s\right)^{2} d t}+\Lambda_{2}(\lambda, \bar{\lambda}) \sqrt{\int_{a}^{\infty}\left(\int_{a}^{\infty} k_{2}(s, t)\right)^{2} d t}
$$

This inequality implies the continuity of the operator $S$.

If conditions I.b') and I.c') are satisfied we can use the fiber Picard theorem again. Consider the spaces $\mathrm{V}=\mathrm{W}=\mathrm{L}^{2}(\mathrm{I})$ and the operators $\mathrm{B}: \mathrm{V} \rightarrow \mathrm{V}$, $\mathrm{C}: \mathrm{V} \times \mathrm{W} \rightarrow \mathrm{W}$ defined by the following relations

$$
\mathrm{B}[v](\mathrm{t})=\mathrm{f}(\mathrm{t})+\int_{\mathrm{a}}^{\mathrm{t}} \mathrm{~K}_{1}(\mathrm{t}, \mathrm{~s}, \lambda, \mathrm{y}(\mathrm{~s})) \mathrm{d} s+\int_{\mathrm{a}}^{\infty} \mathrm{K}_{2}(\mathrm{t}, \mathrm{~s}, \lambda, \mathrm{y}(\mathrm{~s})) \mathrm{ds}
$$

and

$$
\begin{aligned}
C[(v, w)](\mathrm{t})= & \int_{a}^{\mathrm{t}} \frac{\partial \mathrm{~K}_{1}(\mathrm{t}, \mathrm{~s}, v(\mathrm{~s}) ; \lambda)}{\partial \lambda} \mathrm{d} s+\int_{a}^{\infty} \frac{\partial K_{2}(\mathrm{t}, \mathrm{~s}, v(\mathrm{~s}) ; \lambda)}{\partial \lambda} \mathrm{d} s \\
& +\int_{a}^{\mathrm{t}} \frac{\partial \mathrm{~K}_{1}(\mathrm{t}, \mathrm{~s}, v(\mathrm{~s}) ; \lambda)}{\partial v} w(s) \mathrm{d} s+\int_{a}^{\infty} \frac{\partial K_{2}(\mathrm{t}, \mathrm{~s}, v(s) ; \lambda)}{\partial v} w(s) \mathrm{ds}
\end{aligned}
$$

Due to the given conditions B is a Picard operator (condition I.b') implies condition III.) and C satisfies the uniform contraction condition:

$$
\left\|\mathrm{C}\left[\left(v, w_{1}\right)\right]-\mathrm{C}\left[\left(v, w_{2}\right)\right]\right\|_{\mathrm{L}^{2}(\mathrm{I})} \leq \mathrm{L}_{1}\left\|w_{1}-w_{2}\right\|_{\mathrm{L}^{2}(\mathrm{I})}
$$

where $L_{1}=\sqrt{\int_{a}^{\infty} \int_{a}^{t}\left(k_{1}(t, s)+k_{2}(t, s)\right)^{2} d s d t+\int_{a}^{\infty} \int_{t}^{\infty} k_{2}^{2}(t, s) d s d t \text {. Theorem } 1}$ implies that the triangular operator $\mathrm{A}[v, w]=(\mathrm{B}[v], \mathrm{C}[v, w])$ is a Picard operator. Hence the sequence of successive approximation $\left(y_{n+1}, z_{n+1}\right)=A\left[y_{n}, z_{n}\right]$ converges in $\left(\mathrm{L}^{2}(\mathrm{I})\right)^{2}$ to the unique fixed point. If we choose $\mathrm{y}_{\mathrm{o}}(\cdot, \lambda)$ continuously differentiable (with respect to $\lambda$ ) and $z_{0}=\frac{\partial y_{0}}{\partial \lambda}$, then from the construction of the operator $C$ we obtain $z_{n}=\frac{\partial y_{n}}{\partial \lambda}$. On the other hand the operators $S_{n}:\left[\lambda_{1}, \lambda_{2}\right] \rightarrow L^{2}(I)$ defined by $S_{n}(\lambda)(t)=y_{n}(t, \lambda), \forall t \in I, \forall \lambda \in\left[\lambda_{1}, \lambda_{2}\right]$ are differentiables and the differential of $S_{n}$ in $\lambda$ is $z_{n}$, so from Theorem 2 we obtain the differentiability of $S$.

## Acknowledgements

The author wishes to express his deep sense of gratitude to the unknown referee for his valuable observations. This research was supported by a Grant of the Romanian National Authority for Scientific Research, CNCS UEFISCDI, 393 Project Number PN-II-ID-PCE-2011-3-0094.

## References

[1] Sz. András, Fredholm-Volterra equations, Pure Math. Appl. (PU.M.A.), 13 (2002):1-2, 21-30.
[2] Sz. András, Fiber Picard operators and convex contractions, Fixed Point Theory, 4 (2003):2, 209-217.
[3] Sz. András, Fiber $\varphi$-contractions on generalized metric spaces and application, Mathematica (Cluj-Napoca), 45 (2003):1, 3-8.
[4] V. Lakshmikantham, D. Guo, X. Liu, Nonlinear integral equations in abstract spaces, Kluwer Academic Publishers, Dordrecht-Boston-London, 1996.
[5] B. G. Pachpatte, On the existence and uniqueness of solutions of VolterraFredholm integral equations, Mathematics Seminar Notes, 10 (1982), 733-742.
[6] A. Petruşel, Fredholm-Volterra integral equations and Maia's theorem, Preprint, 1988, No. 3, 79-82. Universitatea Babeş-Bolyai.
[7] I. A. Rus, Picard operators and applications, Preprint No. 3, Universitatea Babeş-Bolyai, Cluj-Napoca, 1996.
[8] I. A. Rus, Fiber Picard operators and applications, Mathematica (ClujNapoca), 1999.
[9] I. A. Rus, Fiber Picard operators on generalized metric spaces and application, Scripta Sc. Math., 1 (1999):2, 355-363.
[10] I. A. Rus, Picard operators and applications, Sci. Math. Jpn., 58 (2003):1, 191-219.
[11] M. A. Şerban, Fiber $\varphi$-contractions, Stud. Univ. Babeş-Bolyai Math., 44 (1999):3, 99-108.

# Some inequalities in bicentric quadrilateral 

Mihály Bencze<br>505600 Săcele-Négyfalu,<br>Jud. Braşov, Romania<br>email: benczemihaly@gmail.com

Marius Drăgan<br>061311 bd. Timişoara nr. 35,<br>bl. OD6, sc. E, et. 7 ap. 176, Sect. 6., Bucureşti, Romania<br>email: marius.dragan2005@yahoo.com

## Dedicated to the memory of Professor Antal Bege


#### Abstract

In this paper we prove some results concerning bicentric quadrilaterals. We offer a new proof of the Blundon-Eddy inequality, which we use to obtain other inequalities in bicentric quadrilaterals.


## 1 Introduction

Let $A B C D$ be a bicentric quadrilateral with $a=A B, b=B C, c=C D, d=$ $A D, d_{1}=A C, d_{2}=B D, s=\frac{a+b+c+d}{2}, R$ the radius of the circumscribed circle of the quadrilateral $A B C D$ and $r$ the radius of the inscribed circle, $F$ the area.

In [1] W. J. Blundon and R. H. Eddy proved that:

$$
8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right) \leq s^{2} \leq(r+\sqrt{4 R+r})^{2}
$$

In the following we give a simple proof to this double inequality using the product

$$
(a-b)^{2}(a-c)^{2}(a-d)^{2}(b-c)^{2}(b-d)^{2}(c-d)^{2}
$$

then we deduce many other important new inequalities. We mention that the result concerning the above product is new.

We denote:

$$
\sigma_{1}=\sum a, \sigma_{2}=\sum a b, \sigma_{3}=\sum a b c, x_{1}=b c+a d, x_{2}=a b+c d, x_{3}=a c+b d
$$

[^0]Key words and phrases: bicentric quadrilateral


## 2 Main results

Lemma 1 In every bicentric quadrilateral ABCD the following equalities are true:

1) $F^{2}=(s-a)(s-b)(s-c)(s-d)=a b c d$;
2) $x_{1} x_{2} x_{3}=16 R^{2} r^{2} s^{2}$;
3) $x_{1}+x_{2}=s^{2}$;
4) $x_{1}+x_{2}+x_{3}=s^{2}+2 r^{2}+2 r \sqrt{r^{2}+4 R^{2}}$;
5) $x_{3}=2 r\left(r+\sqrt{4 R^{2}+r^{2}}\right)$;
6) $(a-b)^{2}(a-c)^{2}(a-d)^{2}(b-c)^{2}(b-d)^{2}(c-d)^{2}=\left(x_{1}-x_{2}\right)^{2}\left(x_{2}-x_{3}\right)^{2}$ $\left(x_{3}-x_{1}\right)^{2}$.

## Proof.

1) We have $a+c=b+d$. It results that $s-b=d$ and three similar equalities which imply

$$
(s-a)(s-b)(s-c)(s-d)=a b c d
$$

2) From Ptolemy's theorem it results that $x_{3}=d_{1} d_{2}$. We have the equalities:

$$
a d \sin A+b c \sin C=2 F, a b \sin B+d c \sin D=2 F
$$

We obtain $(a d+b c) d_{1}=4 R F,(a b+d c) d_{2}=4 R F$ which implies

$$
\begin{equation*}
(a d+b c)(a b+d c) d_{1} d_{2}=16 R^{2} F^{2} \text { or } x_{1} x_{2} x_{3}=16 R^{2} r^{2} s^{2} \tag{1}
\end{equation*}
$$

3) We have $x_{1}+x_{2}=a d+b c+a b+c d=(a+c)(d+b)=(a+c)^{2}=$ $\left(\frac{a+b+c+d}{2}\right)=s^{2}$.
4) From (1) it results that

$$
\begin{align*}
& (a b+b c)(a d+d c)(a c+b d)=16 R^{2} F^{2} \text { or } \\
& a b c d \sum a^{2}+\sigma_{3}^{2}-2 a b c d \sigma_{2}=16 R^{2} F^{2} \text { or }  \tag{2}\\
& \sigma_{3}^{2}-4 s^{2} r^{2} \sigma_{2}+4 s^{4} r^{2}=16 R^{2} r^{2} s^{2} v
\end{align*}
$$

But $(s-a)(s-b)(s-c)(s-d)=s^{2} r^{2}$ or $-s^{3}+\sigma_{2} s-\sigma_{3}=0$ which implies

$$
\begin{equation*}
\sigma_{3}=s\left(\sigma_{2}-s^{2}\right) \tag{3}
\end{equation*}
$$

From (2) and (3) we have:

$$
\begin{aligned}
& s^{2}\left(\sigma_{2}-s^{2}\right)^{2}-4 s^{2} r^{2} \sigma_{2}+4 s^{4} r^{2}=16 R^{2} r^{2} s^{2} \quad \text { or } \\
& \sigma_{2}^{2}-\left(2 s^{2}+4 r^{2}\right) \sigma_{2}+s^{4}+2 s^{2} r^{2}-16 r^{2} R^{2}=0
\end{aligned}
$$

It results that: $\sigma_{2}=s^{2}+2 r^{2}+2 r \sqrt{r^{2}+4 R^{2}}$. But $\sigma_{2}=x_{1}+x_{2}+x_{3}$, so it follows that

$$
\begin{equation*}
x_{1}+x_{2}+x_{3}=s^{2}+2 r^{2}+2 r \sqrt{r^{2}+4 R^{2}} \tag{4}
\end{equation*}
$$

5) From 4) since $x_{1}+x_{2}=s^{2}$ it follows that $x_{3}=2 r^{2}+2 r \sqrt{4 R^{2}+r^{2}}$.
6) We have $(a-b)^{2}(a-c)^{2}(a-d)^{2}(b-c)^{2}(b-d)^{2}(c-d)^{2}=$ $[(a-b)(c-d)]^{2}[(a-c)(b-d)]^{2}[(a-d)(b-c)]^{2}=$ $\left(x_{1}-x_{2}\right)^{2}\left(x_{2}-x_{3}\right)^{2}\left(x_{2}-x_{1}\right)^{2}$.

Theorem 1 In every bicentric quadrilateral $\operatorname{ABCD}$ the following equality is true:

$$
\begin{aligned}
& (a-b)^{2}(a-c)^{2}(a-d)^{2}(b-c)^{2}(b-d)^{2}(c-d)^{2} \\
& =16 r^{4} s^{2}\left[s^{2}-8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)\right]\left[s^{2}-\left(r+\sqrt{4 R^{2}+r^{2}}\right)^{2}\right]^{2}
\end{aligned}
$$

Proof. We denote $\triangle=(a-b)^{2}(a-c)^{2}(a-d)^{2}(b-c)^{2}(b-d)^{2}(c-d)^{2}$. From Lemma 1 6) we have:

$$
\begin{align*}
\Delta & =\left(x_{1}-x_{2}\right)^{2}\left(x_{3}-x_{1}\right)^{2}\left(x_{3}-x_{2}\right)^{2} \\
& =\left[\left(x_{1}+x_{2}\right)^{2}-4 x_{1} x_{2}\right]\left[x_{3}^{2}-x_{3}\left(x_{1}+x_{2}\right)+x_{1} x_{2}\right]^{2} . \tag{5}
\end{align*}
$$

From Lemma 1 2) and 5) it results that:

$$
\begin{equation*}
x_{1} x_{2}=\frac{8 R^{2} r^{2} s^{2}}{r\left(r+\sqrt{4 R^{2}+r^{2}}\right)}=2 r\left(\sqrt{4 R^{2}+r^{2}}-r\right) s^{2} \tag{6}
\end{equation*}
$$

From Lemma 1 3), 5) and equalities (5), (6) we obtain:

$$
\begin{aligned}
\triangle= & {\left[s^{4}-8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right) s^{2}\right]\left[4 r^{2}\left(r+\sqrt{4 R^{2}+r^{2}}\right)^{2}\right.} \\
& \left.-2 s^{2} r\left(r+\sqrt{4 R^{2}+r^{2}}\right)+2 r\left(\sqrt{4 R^{2}+r^{2}}-r\right) s^{2}\right]^{2} \\
= & s^{2}\left[s^{2}-8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)\right]\left[4 r^{2}\left(r+\sqrt{4 R^{2}+r^{2}}\right)^{2}-4 r^{2} s^{2}\right]^{2} \\
= & 16 r^{4} s^{2}\left[s^{2}-8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)\right]\left[s^{2}-\left(r+\sqrt{4 R^{2}+r^{2}}\right)^{2}\right] .
\end{aligned}
$$

Theorem 2 In every bicentric quadrilateral ABCD the following double inequality is true: $8 \mathrm{r}\left(\sqrt{4 \mathrm{R}^{2}+\mathrm{r}^{2}}-\mathrm{r}\right) \leq \mathrm{s}^{2} \leq\left(\mathrm{r}+\sqrt{4 \mathrm{R}^{2}+\mathrm{r}^{2}}\right)^{2}$. The equality holds in the case of two bicentric quadrilaterals $\mathrm{A}_{1} \mathrm{~B}_{1} \mathrm{C}_{1} \mathrm{D}_{1}$ and $\mathrm{A}_{2} \mathrm{~B}_{2} \mathrm{C}_{2} \mathrm{D}_{2}$ with the sides

$$
\begin{aligned}
& a_{1}=c_{1}=\sqrt{2 r \sqrt{4 R^{2}+r^{2}}-2 r^{2}} \\
& b_{1}=\sqrt{2 r \sqrt{4 R^{2}+r^{2}}-2 r^{2}}-\sqrt{2 r \sqrt{4 R^{2}+r^{2}}-6 r^{2}} \\
& d_{1}=\sqrt{2 r \sqrt{4 R^{2}+r^{2}}-2 r^{2}}+\sqrt{2 r \sqrt{4 R^{2}+r^{2}}-6 r^{2}} \\
& a_{2}=d_{2}=\frac{r+\sqrt{r^{2}+4 R^{2}}-\sqrt{4 R^{2}-2 r^{2}-2 r \sqrt{4 R^{2}+r^{2}}}}{2} \\
& b_{2}=c_{2}=\frac{r+\sqrt{r^{2}+4 R^{2}}+\sqrt{4 R^{2}-2 r^{2}-2 r \sqrt{4 R^{2}+r^{2}}}}{2} .
\end{aligned}
$$

Proof. We have $\left(x_{3}-x_{1}\right)\left(x_{3}-x_{2}\right)=(a-b)(b-c)(c-d)(d-a)$ and because $a+c=b+d$ it results that $(a-b)(b-c)(c-d)(d-a)=(a-b)^{2}$ $(b-c)^{2} \geq 0$, which implies $\left(x_{3}-x_{1}\right)\left(x_{3}-x_{2}\right) \geq 0$ or

$$
s^{2} \leq\left(r+\sqrt{4 R^{2}+r^{2}}\right)^{2}
$$

But, from Theorem 1 since $\triangle \geq 0$, it results that

$$
8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right) \leq s^{2}
$$

It remain to study the equality cases for $s_{1} \leq s \leq s_{2}$ where

$$
s_{1}=\sqrt{8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)}, s_{2}=r+\sqrt{4 R^{2}+r^{2}}
$$

From Theorem 1 it results that we may have the cases:

## Case 1.

$$
a=c .
$$

We denote $a=x$. Then

$$
a=x, b=y, c=x, d=2 x-y .
$$

From Lemma 1 we have:

$$
x_{3}=2 r\left(r+\sqrt{4 R^{2}+r^{2}}\right) \text { or } x^{2}+y(2 x-y)=2 r\left(r+\sqrt{4 R^{2}+r^{2}}\right) .
$$

But $F^{2}=$ abcd or $(2 x-y) y=4 r^{2}$. It results that $x^{2}=2 r \sqrt{4 R^{2}+r^{2}}-2 r^{2}$. Since $s_{1}^{2}=4 x^{2}=8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)$ represents the left side of the inequality from the statement, so:

$$
\begin{aligned}
& x=\sqrt{2 r \sqrt{4 R^{2}+r^{2}}-2 r^{2}} \\
& (y-x)^{2}=2 r \sqrt{4 R^{2}+r^{2}}-6 r^{2} \text { or }|y-x|=\sqrt{2 r \sqrt{4 R^{2}+r^{2}}-6 r^{2}} .
\end{aligned}
$$

We denote $u_{1}=2 r \sqrt{4 R^{2}+r^{2}}-2 r^{2}, u_{2}=2 r \sqrt{4 R^{2}+r^{2}}-6 r^{2}$. If $x \leq y$ we have

$$
a=x=\sqrt{u_{1}}, b=y=\sqrt{u_{1}}+\sqrt{u_{2}}, c=\sqrt{u_{1}}, d=2 x-y=\sqrt{u_{1}}-\sqrt{u_{2}} .
$$

If $x>y$ we have
$a=x=\sqrt{u_{1}}, b=y=x-\sqrt{u_{2}}=\sqrt{u_{1}}-\sqrt{u_{2}}, c=\sqrt{u_{1}}, d=2 x-y=\sqrt{u_{1}}+\sqrt{u_{2}}$.
It results that the equality from the left side of the inequality of the statement holds in the case of bicentric quadrilateral $A_{1} B_{1} C_{1} D_{1}$ with the sides

$$
\sqrt{u_{1}}, \sqrt{u_{1}}-\sqrt{u_{2}}, \sqrt{u_{1}}, \sqrt{u_{1}}+\sqrt{u_{2}} .
$$

## Case 2.

$$
a=d=x, b=c=y .
$$

In this case $m(\measuredangle D)=m(\measuredangle B)=90^{\circ}, A C=2 R$. It results that $F=s r=2 \frac{x y}{2}$ or $x y=(x+y)$ r.

We denote $\alpha=x+y, \beta=x y$.
We have $\beta=\alpha$. But $x^{2}+y^{2}=4 R^{2}$ which implies $\alpha^{2}-2 \beta=4 R^{2}$ so we have $\alpha^{2}-2 \alpha r-4 R^{2}=0$.

It results that $\alpha=r+\sqrt{r^{2}+4 R^{2}}$.
But $s_{1}=x+y=\alpha=r+\sqrt{r^{2}+4 R^{2}}$ which represents the right side of the inequality from the statement. We have $\left\{\begin{array}{l}x+y=\alpha \\ x y=r \alpha\end{array}\right.$, so $x, y$ are the solutions of the equation $\mathfrak{u}^{2}-\alpha u+r \alpha=0$ which implies:

$$
\begin{aligned}
& x=\frac{\alpha-\sqrt{\alpha^{2}-4 r \alpha}}{2}=\frac{r+\sqrt{r^{2}+4 R^{2}}-\sqrt{4 R^{2}-2 r^{2}-2 r \sqrt{4 R^{2}+r^{2}}}}{2} \\
& y=\frac{r+\sqrt{r^{2}+4 R^{2}}+\sqrt{4 R^{2}-2 r^{2}-2 r \sqrt{4 R^{2}+r^{2}}}}{2}
\end{aligned}
$$

So, the equality for the right side of the inequality from the statement is true in the case of bicentric quadrilateral $\mathrm{A}_{2} \mathrm{~B}_{2} \mathrm{C}_{2} \mathrm{D}_{2}$ with the sides

$$
a_{2}=x, b_{2}=x, c_{2}=y, d_{2}=y .
$$

Theorem 3 In every bicentric quadrilateral ABCD the following inequalities are true:

$$
\begin{aligned}
2 r\left(r+\sqrt{4 R^{2}+r^{2}}\right) & \leq \min \{a b+c d, b c+a d\} \leq 4 r\left(\sqrt{4 R^{2}+r^{2}}-r\right) \\
& \leq \max \{a b+c d+b c+a d\} \leq 4 R^{2} .
\end{aligned}
$$

Proof. We suppose that $x_{1} \leq x_{2}, x_{1}+x_{2}=s^{2}, x_{1} x_{2}=\alpha s^{2}$ where

$$
\alpha=\frac{8 R^{2} r}{\sqrt{4 R^{2}+r^{2}}+r}=2 r\left(\sqrt{4 R^{2}+r^{2}}-r\right) .
$$

It results that: $x_{1}=\frac{s^{2}-\sqrt{s^{4}-4 \alpha s^{2}}}{2}, x_{2}=\frac{s^{2}+\sqrt{s^{4}-4 \alpha s^{2}}}{2}$. We consider the functions $\mathrm{f}, \mathrm{g}:(0,+\infty) \rightarrow \mathrm{R}$.

$$
f(s)=\frac{s^{2}-\sqrt{s^{4}-4 \alpha s^{2}}}{2}, g(s)=\frac{s^{2}+\sqrt{s^{4}-4 \alpha s^{2}}}{2} .
$$

After differentiation we obtain:

$$
f^{\prime}(s)=\frac{s\left(\sqrt{s^{4}-4 \alpha s^{2}}-s^{2}+2 \alpha\right)}{\sqrt{s^{4}-4 \alpha s^{2}}} \leq 0, g^{\prime}(s)=\frac{s\left(\sqrt{s^{4}-4 \alpha s^{2}}+s^{2}-4 \alpha\right)}{\sqrt{s^{4}-4 \alpha s^{2}}} \geq 0
$$

From Theorem 2 it results that: $s^{2} \geq 8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)=4 \alpha$.
It results that f is a decreasing and g is an increasing function. Because $s \leq r+\sqrt{4 R^{2}+r^{2}}$ we have $f\left(r+\sqrt{4 R^{2}+r^{2}}\right) \leq f(s)=x_{1}$. If follows that

$$
\begin{aligned}
x_{1} \geq & \frac{1}{2}\left[\left(r+\sqrt{4 R^{2}+r^{2}}\right)^{2}\right. \\
& \left.-\left(r+\sqrt{4 R^{2}+r^{2}}\right) \sqrt{\left(r+\sqrt{4 R^{2}+r^{2}}\right)^{2}-8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)}\right] \\
= & \frac{\left(r+\sqrt{4 R^{2}+r^{2}}\right)}{2}\left[r+\sqrt{4 R^{2}+r^{2}}\right. \\
& \left.-\sqrt{r^{2}+4 R^{2}+r^{2}+2 r \sqrt{4 R^{2}+r^{2}}-8 r \sqrt{4 R^{2}+r^{2}}+8 r^{2}}\right] \\
= & \frac{\left(r+\sqrt{4 R^{2}+r^{2}}\right)}{2}\left[r+\sqrt{4 R^{2}+r^{2}}-\sqrt{\left(\sqrt{4 R^{2}+r^{2}}\right)^{2}+9 r^{2}-6 r \sqrt{4 R^{2}+r^{2}}}\right] \\
= & 2 r\left(r+\sqrt{4 R^{2}+r^{2}}\right) .
\end{aligned}
$$

It follows that

$$
\begin{equation*}
x_{1} \geq 2 r\left(r+\sqrt{4 R^{2}+r^{2}}\right) \tag{7}
\end{equation*}
$$

From $s \leq r+\sqrt{4 R^{2}+r^{2}}$ it results also that

$$
\begin{aligned}
x_{2}= & g(s) \leq g\left(r+\sqrt{4 R^{2}+r^{2}}\right) \\
= & \frac{1}{2}\left[\left(r+\sqrt{4 R^{2}+r^{2}}\right)^{2}\right. \\
& \left.+\left(r+\sqrt{4 R^{2}+r^{2}}\right) \sqrt{\left(r+\sqrt{4 R^{2}+r^{2}}\right)^{2}-8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)}\right] \\
= & \left(\sqrt{4 R^{2}+r^{2}}+r\right)\left(\sqrt{4 R^{2}+r^{2}}-r\right)=4 R^{2} .
\end{aligned}
$$

Thus we get the following inequality

$$
\begin{equation*}
x_{2} \leq 4 R^{2} \tag{8}
\end{equation*}
$$

Since $8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right) \leq s^{2}$ we have $x_{1}=f(s) \leq f\left(\sqrt{8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)}\right)$ or in an equivalent form

$$
\begin{aligned}
x_{1} \leq & \frac{1}{2}\left[8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)\right. \\
& \left.-\sqrt{8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)} \sqrt{8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)-8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)}\right] \\
= & 4 r\left(\sqrt{4 R^{2}+r^{2}}-r\right) .
\end{aligned}
$$

It follows that

$$
\begin{equation*}
x_{1} \leq 4 r\left(\sqrt{4 R^{2}+r^{2}}-r\right) \tag{9}
\end{equation*}
$$

Because $8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right) \leq s^{2}$ and $g$ is an increasing function it follows that:

$$
\begin{equation*}
g\left(\sqrt{8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)}\right) \leq g(s)=x_{2} \text { or } x_{2} \geq 4 r\left(\sqrt{4 R^{2}+r^{2}}-r\right) \tag{10}
\end{equation*}
$$

From (7) (8) (9) and (10) it results that:

$$
x_{3}=2 r\left(r+\sqrt{4 R^{2}+r^{2}}\right) \leq x_{1} \leq 4 r\left(\sqrt{4 R^{2}+r^{2}}-r\right) \leq x_{2} \leq 4 R^{2}
$$

Remark 1 From Theorem 3 it results that $2 r\left(r+\sqrt{4 R^{2}+r^{2}}\right) \leq 4 r\left(\sqrt{4 R^{2}+r^{2}}-\right.$ r) which, after performing some calculation, represent the well-known Fejes inequality $R \geq \sqrt{2} r$.

Theorem 4 In every bicentric quadrilateral ABCD the following inequalities are true:

$$
\begin{aligned}
\frac{r\left(\sqrt{4 R^{2}+r^{2}}+r\right)}{R} & \leq \min \left\{d_{1}, d_{2}\right\} \leq \frac{\sqrt{4 R^{2}+r^{2}}+r}{R} \sqrt{\frac{\left(\sqrt{4 R^{2}+r^{2}}-r\right) r}{2}} \\
& \leq \max \left\{d_{1}, d_{2}\right\} \leq 2 R
\end{aligned}
$$

Proof. We suppose that $x_{1} \leq x_{2}$.
From Ptolemy's theorem it results that $\frac{x_{1}}{x_{2}}=\frac{d_{1}}{d_{2}}$ which implies $d_{1} \leq d_{2}$.
Because $d_{1} d_{2}=x_{3}$ we have

$$
\begin{aligned}
d_{1}^{2} & =\frac{x_{1}}{x_{2}} x_{3}=\frac{s^{2}-\sqrt{s^{4}-4 \alpha s^{2}}}{s^{2}+\sqrt{s^{4}-4 \alpha s^{2}}} x_{3}=x_{3} \frac{\left(s^{2}-\sqrt{s^{4}-4 \alpha s^{2}}\right)^{2}}{4 \alpha s^{2}} \\
& =x_{3} \frac{2 s^{4}-4 \alpha s^{2}-2 s^{2} \sqrt{s^{4}-4 \alpha s^{2}}}{4 \alpha s^{2}}=\frac{x_{3}\left(s^{2}-2 \alpha-\sqrt{s^{4}-4 \alpha s^{2}}\right)}{2 \alpha} \\
& =\frac{2 r\left(r+\sqrt{4 R^{2}+r^{2}}\right)}{4 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)}\left[s^{2}-\sqrt{s^{4}-4 \alpha s^{2}}-2 \alpha\right] \\
& =\frac{\left(\sqrt{4 R^{2}+r^{2}}+r\right)^{2}}{8 R^{2}}\left[s^{2}-\sqrt{s^{4}-4 \alpha s^{2}}-2 \alpha\right]=B\left(2 x_{1}-2 \alpha\right)
\end{aligned}
$$

where we denote $B=\frac{\left(\sqrt{4 R^{2}+r^{2}}+r\right)^{2}}{8 R^{2}}$.
But from Theorem 3 we have

$$
4 r\left(r+\sqrt{4 R^{2}+r^{2}}\right) \leq 2 x_{1} \leq 8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)
$$

We obtain

$$
\begin{aligned}
& 4 r\left(r+\sqrt{4 R^{2}+r^{2}}\right)-2 \alpha \leq 2 x_{1}-2 \alpha \leq 8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)-2 \alpha \text { or } \\
& 8 r^{2} \leq 2 x_{1}-2 \alpha \leq 4 r\left(\sqrt{4 R^{2}+r^{2}}-r\right) \text { or } \\
& 8 r^{2} B \leq B\left(2 x_{1}-2 \alpha\right) \leq 4 r\left(\sqrt{4 R^{2}+r^{2}}-r\right) B \text { or } \\
& \frac{8 r^{2}\left(\sqrt{4 R^{2}+r^{2}}+r\right)^{2}}{8 R^{2}} \leq d_{1}^{2} \leq \frac{4 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)\left(\sqrt{4 R^{2}+r^{2}}+r\right)^{2}}{8 R^{2}}
\end{aligned}
$$

It results that:

$$
\begin{equation*}
\frac{r\left(\sqrt{4 R^{2}+r^{2}}+r\right)}{r}<d_{1} \leq \frac{\sqrt{4 R^{2}+r^{2}}+r}{R} \sqrt{\frac{\left(\sqrt{4 R^{2}+r^{2}}-r\right) r}{2}} \tag{11}
\end{equation*}
$$

Also:

$$
\begin{aligned}
d_{2}^{2} & =\frac{x_{2}}{x_{1}} x_{3}=\frac{s^{2}+\sqrt{s^{4}-4 \alpha s^{2}}}{s^{2}-\sqrt{s^{4}-4 \alpha s^{2}}} x_{3}=\frac{\left(s^{2}+\sqrt{s^{4}-4 \alpha s^{2}}\right)^{2}}{4 \alpha s^{2}} x_{3} \\
& =\frac{x_{3}}{4 \alpha s^{2}}\left(2 s^{4}-4 \alpha s^{2}+2 s^{2} \sqrt{s^{4}-4 \alpha s^{2}}\right)=\frac{x_{3}}{2 \alpha}\left(s^{2}+\sqrt{s^{4}-4 \alpha s^{2}}-2 \alpha\right) \\
& =\frac{x_{3}}{2 \alpha}\left(2 x_{2}-2 \alpha\right)=\frac{\left(\sqrt{4 R^{2}+r^{2}}+r\right)^{2}\left(2 x_{2}-2 \alpha\right)}{8 R^{2}}
\end{aligned}
$$

But we have proved that $4 r\left(\sqrt{4 R^{2}+r^{2}}-r\right) \leq x_{2} \leq 4 R$. It results that:

$$
\begin{align*}
& 4 r\left(\sqrt{4 R^{2}+r^{2}}-r\right) \leq 2 x_{2}-2 \alpha \leq 2\left(4 R^{2}+2 r^{2}-2 r \sqrt{4 R^{2}+r^{2}}\right) \text { or } \\
& 4 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)\left(\frac{\sqrt{4 R^{2}+r^{2}}+r}{2 \sqrt{2} R}\right)^{2} \leq d_{2}^{2} \\
& \quad \leq 2\left(\sqrt{4 R^{2}+r^{2}}-r\right)^{2} \frac{\left(\sqrt{4 R^{2}+r^{2}}+r\right)^{2}}{8 R^{2}} \text { or }  \tag{12}\\
& \frac{\sqrt{4 R^{2}+r^{2}}+r}{R} \sqrt{\frac{\left(\sqrt{4 R^{2}+r^{2}}-r\right) r}{2} \leq d_{2} \leq 2 R .}
\end{align*}
$$

From (11) and (12) it results the inequalities from the statement.
Theorem 5 Let be $\alpha, \beta \in R$ so that $s \leq \alpha R+\beta r$ is true in every bicentric quadrilateral ABCD . Then $2 \mathrm{R}+(4-2 \sqrt{2}) \mathrm{r} \leq \alpha \mathrm{R}+\beta \mathrm{r}$ is true in every bicentric quadrilateral ABCD.

Proof. We consider the case of the square with the sides $\mathrm{a}=\mathrm{b}=\mathrm{c}=\mathrm{d}=1$. We have $2 \leq \alpha \frac{1}{\sqrt{2}}+\beta \frac{1}{2}$. It results that

$$
\begin{equation*}
4 \leq \sqrt{2} \alpha+\beta \tag{13}
\end{equation*}
$$

If $a=b=1, c=d=0$ it results that $R=\frac{1}{2}, r=0$.
It follows that

$$
\begin{equation*}
1 \leq \frac{\alpha}{2} \text { or } \alpha \geq 2 \tag{14}
\end{equation*}
$$

We know that

$$
\begin{equation*}
R \geq \sqrt{2} r \tag{15}
\end{equation*}
$$

From (13), (14) and (15) it results that

$$
\begin{aligned}
(\alpha-2) R+(\beta-4+2 \sqrt{2}) r & \geq(\alpha-2) \sqrt{2} r+(\beta-4+2 \sqrt{2}) r \\
& =(\alpha \sqrt{2}+\beta-4) r \geq 0
\end{aligned}
$$

therefore

$$
\alpha R+\beta r \geq 2 R+(4-2 \sqrt{2}) r
$$

Theorem 6 In every bicentric quadrilateral the following inequality is true:

$$
s \leq 2 R+(4-2 \sqrt{2}) r
$$

Proof. From the Theorem 1 we have $s \leq r+\sqrt{4 R^{2}+r^{2}}$. We denote $x=\frac{R}{r}$.
We prove that

$$
r+\sqrt{4 R^{2}+r^{2}} \leq 2 R+(4-2 \sqrt{2}) r
$$

or in an equivalent form

$$
\begin{aligned}
& 1+\sqrt{4 x^{2}+1} \leq 2 x+4-2 \sqrt{2} \text { or } \sqrt{4 x^{2}+1} \leq 2 x+3-2 \sqrt{2} \text { or } \\
& 1 \leq 4(3-2 \sqrt{2}) x+(3-2 \sqrt{2})^{2} \text { or } x \geq \frac{(-2+2 \sqrt{2})(4-2 \sqrt{2})}{4(3-2 \sqrt{2})}
\end{aligned}
$$

After performing some calculation it results that $x \geq \sqrt{2}$ which represents just the Fejes's inequality [2].

Theorem 7 In every bicentric quadrilateral ABCD the following inequalities are true:

1) $4 r\left(3 \sqrt{4 R^{2}+r^{2}}-5 r\right) \leq a^{2}+b^{2}+c^{2}+d^{2} \leq 8 R^{2}$;
2) $2 r \sqrt{8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)}\left(7 \sqrt{4 R^{2}+r^{2}}-9 r\right) \leq \sum a^{2} b \leq 8 R^{2}+2 r^{2}$;
3) $2 r\left(5 \sqrt{4 R^{2}+r^{2}}-3 r\right) \leq \sum a b \leq 4\left(R^{2}+r^{2}+r \sqrt{4 R^{2}+r^{2}}\right)$;
4) $32 r^{2} \sqrt{4 R^{2}+r^{2}}\left(\sqrt{4 R^{2}+r^{2}}-r\right) \leq \sum a^{2} b c$
$\leq 4 r \sqrt{4 R^{2}+r^{2}}\left(r+\sqrt{4 R^{2}+r^{2}}\right)^{2} ;$
5) $\left(2 r^{2}+2 r \sqrt{4 R^{2}+r^{2}}\right) \sqrt{8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)} \leq \sum a b c$
$\leq 2 r\left(r+\sqrt{4 R^{2}+r^{2}}\right)^{2}$.
Proof. We have $\sigma_{2}=s^{2}+\alpha, \sigma_{3}=s \alpha$ where $\alpha=2 r^{2}+2 r \sqrt{r^{2}+4 R^{2}}$.
6) $\sum a^{2}=(2 s)^{2}-2 \sigma_{2}=4 s^{2}-2 \sigma_{2}=4 s^{2}-2 s^{2}-4 r^{2}-4 r \sqrt{4 R^{2}+r^{2}}$.

It results that: $\sum a^{2}=2 s^{2}-4 r^{2}-4 r \sqrt{4 R^{2}+r^{2}}$.
From Theorem 2 we obtain

$$
4 r\left(3 \sqrt{4 R^{2}+r^{2}}-5 r\right) \leq a^{2}+b^{2}+c^{2}+d^{2} \leq 8 R^{2}
$$

2) $a^{2} b=a b(2 s-b-c-d)=2 s a b-a b^{2}-a b c-a b d$ or $a^{2} b+a b^{2}=$ $2 s a b-a b c-a b d$.
It results that $\sum a^{2} b=2 s \sigma_{2}-3 \sigma_{3}=2 s^{3}-s \alpha=s\left(2 s^{2}-\alpha\right)$ which implies $\sum a^{2} b=s\left(2 s^{2}-\alpha\right)$. We consider the increasing function

$$
\begin{aligned}
& f:(0,+\infty) \rightarrow R, f(s)=2 s^{3}-s \alpha, \text { with } f^{\prime}(s)=6 s^{2}-\alpha \geq 0 \text { as } \\
& s^{2} \geq 8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right) \geq \frac{\alpha}{6}=\frac{2 r^{2}+2 r \sqrt{r^{2}+4 R^{2}}}{6} .
\end{aligned}
$$

The last inequality may be written as:

$$
24 \sqrt{4 R^{2}+r^{2}}-24 r \geq r+\sqrt{4 R^{2}+r^{2}} \text { or } 23 \sqrt{4 R^{2}+r^{2}} \geq 25 r .
$$

But from inequality of Fejes it results that

$$
23 \sqrt{4 R^{2}+r^{2}} \geq 25 \sqrt{9 r^{2}}=75 r>25 r .
$$

From Theorem 2 it results that:

$$
\begin{aligned}
& \sqrt{8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)}\left(16\left(\sqrt{4 R^{2}+r^{2}}-r\right)-2 r^{2}-2 r \sqrt{4 R^{2}+r^{2}}\right) \\
& \leq \sum a^{2} b \leq\left(r+\sqrt{4 R^{2}+r^{2}}\right) \\
& \left(2 r^{2}+8 R^{2}+2 r^{2}+2 r \sqrt{4 R^{2}+r^{2}}-2 r^{2}-2 r \sqrt{4 R^{2}+r^{2}}\right)
\end{aligned}
$$

which is equivalent with the inequality from the statement.
3) $\sigma_{2}=\sum a b=s^{2}+\alpha$ or $8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)+2 r^{2}+2 r \sqrt{4 R^{2}+r^{2}} \leq \sum a b \leq$ $r^{2}+4 R^{2}+r^{2}+2 r \sqrt{4 R^{2}+r^{2}}+2 r^{2}+2 r \sqrt{4 R^{2}+r^{2}}$
which is equivalent with the inequality from the statement.
4) $a^{2} b c=a a b c=(2 s-b-c-d) a b c=2 s a b c-a b^{2} c-a b c^{2}-a b c d$ or
$a^{2} b c+a b^{2} c+a b c^{2}=2 s a b c-a b c d$ or $\sum a^{2} b c=2 s \sigma_{3}-4 a b c d=2 s$
$s \alpha-4 s^{2} r^{2}$ or $\sum a^{2} b c=s^{2}\left(2 \alpha-4 r^{2}\right)=s^{2}\left(4 r^{2}+4 r \sqrt{4 R^{2}+r^{2}}-4 r^{2}\right)=$ $4 r \sqrt{4 R^{2}+r^{2}} s^{2}$.

From Theorem 2 it results the inequality from the statement.
5) $\sum \mathrm{abc}=\mathrm{s} \alpha$.

According to Theorem 2 it results the inequality from the statement.

Theorem 8 Let be $\alpha, \beta, \gamma \in R, \beta \geq 4$ so that $s^{2} \leq \alpha R^{2}+\beta R r+\gamma r^{2}$ is true in all bicentric quadrilateral. Then

$$
4 R^{2}+4 R r+(8-4 \sqrt{2}) r^{2} \leq \alpha R^{2}+\beta R r+\gamma r^{2}
$$

is true in all bicentric quadrilateral.

Proof. We consider the case of the bicentric quadrilateral with $\mathrm{a}=\mathrm{b}=\mathrm{c}=$ $\mathrm{d}=1$. It results that $4 \leq \frac{\alpha}{2}+\frac{\beta}{2 \sqrt{2}}+\frac{\gamma}{4}$ or $16 \leq 2 \alpha+\sqrt{2} \beta+\gamma$.

In the case of $a=b=1, c=d=0$ it results that $R=\frac{1}{2}, r=0$ and $\alpha \geq 4$. But from inequality $R \geq \sqrt{2} r$ we have:

$$
\begin{aligned}
& (\alpha-4) R^{2}+(\beta-4) \operatorname{Rr}+(\gamma-8+4 \sqrt{2}) r^{2} \\
& \geq(\alpha-4) 2 r^{2}+\sqrt{2}(\beta-4) r^{2}+(\gamma-8+4 \sqrt{2}) r^{2} \\
& \geq(\alpha-4) 2 r^{2}+\sqrt{2}(\beta-4) r^{2}+(\gamma-8+4 \sqrt{2}) r^{2} \\
& =(2 \alpha+\sqrt{2} \beta+\gamma-16) r^{2} \geq 0 .
\end{aligned}
$$

Theorem 9 In every bicentric quadrilateral ABCD the following inequality is true:

$$
s^{2} \leq 4 R^{2}+4 R r+(8-4 \sqrt{2}) r^{2}
$$

Proof. Since $s^{2} \leq\left(r+\sqrt{4 R^{2}+r^{2}}\right)^{2}$ it is sufficient to prove that:

$$
\begin{aligned}
& \left(\sqrt{4 x^{2}+1}+1\right)^{2} \leq 4 x^{2}+4 x+8-4 \sqrt{2} \text { or } \\
& 4 x^{2}+1+1+2 \sqrt{4 x^{2}+1} \leq 4 x^{2}+4 x+8-4 \sqrt{2} \text { or } \\
& 2 \sqrt{4 x^{2}+1} \leq 4 x+6-4 \sqrt{2} \text { or } \\
& \sqrt{4 x^{2}+1} \leq 2 x+3-2 \sqrt{2} \text { or } 4 x^{2}+1 \leq 4 x^{2}+(12-8 \sqrt{2}) x+(3-2 \sqrt{2})^{2} \text { or } \\
& x \geq \frac{(1-3+2 \sqrt{2})(1+3-2 \sqrt{2})}{4(3-2 \sqrt{2})}=\frac{(\sqrt{2}-1)(2-\sqrt{2})}{3-2 \sqrt{2}}=\sqrt{2}
\end{aligned}
$$

Theorem 10 In every bicentric quadrilateral ABCD the following inequalities are true:

1) $\sum a b c \leq 8 R^{2} r+8 \operatorname{Rr}^{2}+(16-8 \sqrt{2}) r^{3}$;
2) $\sum a b \leq 4\left[R^{2}+2 R r+(4-2 \sqrt{2}) r^{2}\right]$;
3) $\sum \mathrm{a}^{2} \mathrm{bc} \leq 32 \mathrm{R}^{3} \mathrm{r}+16 \mathrm{Rr}^{3}+(80-32 \sqrt{2}) \mathrm{R}^{2} \mathrm{r}^{2}+(32-16 \sqrt{2}) \mathrm{r}^{4}$.

## Proof.

1) We proved that $\sum a b c \leq 2 r\left(r+\sqrt{4 R^{2}+r^{2}}\right)^{2}$, and

$$
\left(r+\sqrt{4 R^{2}+r^{2}}\right)^{2} \leq 4 R^{2}+4 R r+(8-4 \sqrt{2}) r^{2} .
$$

It results that

$$
\sum a b c \leq 2 r\left(4 R^{2}+4 R r+(8-4 \sqrt{2}) r^{2}\right)
$$

2) Since $\sqrt{4 R^{2}+r^{2}} \leq 2 R+(3-2 \sqrt{2}) r$, from Theorem 7 3) it results that:

$$
\begin{aligned}
\sum a b & \leq 4\left(R^{2}+r^{2}+r \sqrt{4 R^{2}+r^{2}}\right) \\
& \leq 4\left[R^{2}+r^{2}+r(2 R+(3-2 \sqrt{2}) r)\right] \\
& =4\left[R^{2}+r^{2}+2 R r+(3-2 \sqrt{2}) r^{2}\right] \text { or } \\
\sum a b & \leq 4\left[R^{2}+2 R r+(4-2 \sqrt{2}) r^{2}\right] .
\end{aligned}
$$

3) From Theorem 7 4) it results that:

$$
\begin{aligned}
\sum a^{2} b c \leq & 4 r \sqrt{4 R^{2}+r^{2}}\left(r+\sqrt{4 R^{2}+r^{2}}\right)^{2} \\
= & 4 r \sqrt{4 R^{2}+r^{2}}\left(r^{2}+4 R^{2}+r^{2}+2 r \sqrt{4 R^{2}+r^{2}}\right) \\
= & 8 r \sqrt{4 R^{2}+r^{2}}\left(2 R^{2}+r^{2}+r \sqrt{4 R^{2}+r^{2}}\right) \\
= & \left(16 R^{2} r+8 r^{3}\right) \sqrt{4 R^{2}+r^{2}}+8 r^{2}\left(4 R^{2}+r^{2}\right) \\
\leq & \left(16 R^{2} r+8 r^{3}\right)[2 R+(3-2 \sqrt{2}) r]+32 R^{2} r^{2}+8 r^{4} \\
= & 32 R^{3} r+(48-32 \sqrt{2}) R^{2} r^{2}+16 R r^{3}+(24-16 \sqrt{2}) r^{4} \\
& +32 R^{2} r^{2}+8 r^{4},
\end{aligned}
$$

which is equivalent with the inequality from the statement.

Theorem 11 In every bicentric quadrilateral ABCD the following inequalities are true:

1) $2 r \sqrt{8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)}\left(5 \sqrt{4 R^{2}+r^{2}}-11 r\right) \leq \sum a^{3}$
$\leq 2\left(r+\sqrt{4 R^{2}+r^{2}}\right)\left(4 R^{2}-r^{2}-r \sqrt{4 R^{2}+r^{2}}\right) ;$
2) $352 R^{2} r^{2}+208 r^{4}-240 r^{3} \sqrt{4 R^{2}+r^{2}}$
$\leq \sum a^{3} b \leq\left(r+\sqrt{4 R^{2}+r^{2}}\right)^{2}\left(8 R^{2}-4 r^{2}\right)$.

## Proof.

1) $\mathrm{a}^{3}=\mathrm{a}^{2}(2 \mathrm{~s}-\mathrm{b}-\mathrm{c}-\mathrm{d})=2 \mathrm{a}^{2} \mathrm{~s}-\mathrm{a}^{2} \mathrm{~b}-\mathrm{a}^{2} \mathrm{c}-\mathrm{a}^{2} \mathrm{~d}$ or $\sum \mathrm{a}^{3}=2 \mathrm{~s} \sum \mathrm{a}^{2}-$ $\sum a^{2} b=2 s\left(2 s^{2}-2 \alpha\right)-2 s^{3}+s \alpha$.
It results that $\sum a^{3}=2 s^{3}-3 \alpha$.
We consider the function $f:(0,+\infty) \rightarrow R, f(s)=2 s^{3}-3 \alpha s$, with the derivate $f^{\prime}(s)=6 s^{2}-3 \alpha$. We prove that $f^{\prime}(s) \geq 0$ or $s^{2} \geq \frac{\alpha}{2}$.
But $s^{2} \geq 8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)$. It will be sufficient to prove that:

$$
\begin{aligned}
& 8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right) \geq r^{2}+r \sqrt{4 R^{2}+r^{2}} \text { or } \\
& 8 \sqrt{4 x^{2}+1}-8 \geq 1+\sqrt{4 x^{2}+1} \text { or } \sqrt{4 x^{2}+1} \geq \frac{9}{7}
\end{aligned}
$$

which is true because $\sqrt{4 x^{2}+1} \geq 2$ according to Fejes inequality.
Since $f$ is an increasing function it results from Theorem 2 that:

$$
\begin{aligned}
& \sqrt{8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)}\left[16\left(\sqrt{4 R^{2}+r^{2}}-r\right)-6 r^{2}-6 r \sqrt{4 R^{2}+r^{2}}\right] \\
& \leq \sum a^{3} \leq\left(r+\sqrt{4 R^{2}+r^{2}}\right)\left[2 r^{2}+8 R^{2}+2 r^{2}\right. \\
& \left.\quad+4 r \sqrt{4 R^{2}+r^{2}}-6 r^{2}-6 r \sqrt{4 R^{2}+r^{2}}\right],
\end{aligned}
$$

which is equivalent with the inequality from the statement.
2) $\mathrm{a}^{3} \mathrm{~b}=\mathrm{ab}\left(\sum \mathrm{a}^{2}-\mathrm{b}^{2}-\mathrm{c}^{2}-\mathrm{d}^{2}\right)=\mathrm{ab} \sum \mathrm{a}^{2}-\mathrm{ab}^{3}-a b c^{2}-a b d^{2}$ or $a^{3} b+a b^{3}=a b \sum a^{2}-a b c^{2}-a b d^{2}$ or $\sum a^{3} b=\sum a b \sum a^{2}-\sum a^{2} b c=$ $\left(s^{2}+\alpha\right)\left(2 s^{2}-2 \alpha\right)-\left(2 \alpha-4 r^{2}\right) s^{2}$ or $\sum a^{3} b=2 s^{4}-\left(2 \alpha-4 r^{2}\right) s^{2}-2 \alpha^{2}$. We denote $s^{2}=t$ and consider the function: $f:(0,+\infty) \rightarrow R$,

$$
f(t)=2 t^{2}-\left(2 a-4 r^{2}\right) t-2 a^{2}
$$

and

$$
\mathrm{t}_{v}=\frac{2 a-4 \mathrm{r}^{2}}{4}=\frac{a-2 r^{2}}{2}=r \sqrt{4 R^{2}+\mathrm{r}^{2}}
$$

We prove that $t \geq t_{v}$.
$s^{2} \geq r \sqrt{4 R^{2}+r^{2}}$. But $s^{2} \geq 8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right)$. It will be sufficient to prove that

$$
8 r\left(\sqrt{4 R^{2}+r^{2}}-r^{2}\right) \geq r \sqrt{4 R^{2}+r^{2}} \text { or } \sqrt{4 R^{2}+r^{2}} \geq \frac{8}{7}
$$

which is true because $\sqrt{4 R^{2}+r^{2}} \geq 3$.
It results that f is an increasing function which implies:

$$
\begin{aligned}
& 128 r^{2}\left(4 R^{2}+2 r^{2}-2 r \sqrt{4 R^{2}+r^{2}}\right)-4 r \sqrt{4 R^{2}+r^{2}} 8 r\left(\sqrt{4 R^{2}+r^{2}}-r\right) \\
& -2\left(2 r^{2}+2 r \sqrt{4 R^{2}+r^{2}}\right)^{2} \leq \sum a^{3} b \leq 2\left(r+\sqrt{4 R^{2}+r^{2}}\right)^{4} \\
& -4 r \sqrt{4 R^{2}+r^{2}}\left(r+\sqrt{4 R^{2}+r^{2}}\right)^{2}-2\left(2 r^{2}+2 r \sqrt{4 R^{2}+r^{2}}\right)^{2}
\end{aligned}
$$

or

$$
\begin{aligned}
& 512 R^{2} r^{2}+256 r^{4}-256 r^{3} \sqrt{4 R^{2}+r^{2}}-32 r^{2}\left(4 R^{2}+r^{2}\right)+32 r^{3} \sqrt{4 R^{2}+r^{2}} \\
& -8 r^{4}-8 r^{2}\left(4 R^{2}+r^{2}\right)-16 r^{3} \sqrt{4 R^{2}+r^{2}} \leq \sum a^{3} b \leq 2\left(r+\sqrt{4 R^{2}+r^{2}}\right)^{2} \\
& \left(r^{2}+4 R^{2}+r^{2}+2 r \sqrt{4 R^{2}+r^{2}}-2 r \sqrt{4 R^{2}+r^{2}}\right)-8 r^{2}\left(r+\sqrt{4 R^{2}+r^{2}}\right)^{2}
\end{aligned}
$$

or

$$
\begin{aligned}
352 R^{2} r^{2}+208 r^{4}-240 r^{3} \sqrt{4 R^{2}+r^{2}} \leq & \sum a^{3} b \leq\left(r+\sqrt{4 R^{2}+r^{2}}\right)^{2} \\
& \left(4 r^{2}+8 R^{2}-8 r^{2}\right)
\end{aligned}
$$

Theorem 12 In every bicentric quadrilateral ABCD the following inequalities are true:

1) $\sum a^{3} \leq 16 R^{3}+(24-16 \sqrt{2}) R^{2} r-8 R r^{2}-(16-8 \sqrt{2}) r^{3}$;
2) $\sum a^{3} b \leq 32 R^{4}-16 R^{2} r^{2}+32 R^{3} r+16 R r^{3}+(64-32 \sqrt{2}) R^{2} r^{2}-(32-16 \sqrt{2}) r^{4}$;
3) $\sum \mathrm{a}^{3} \mathrm{~b} \geq 352 \mathrm{R}^{2} \mathrm{r}^{2}+(480 \sqrt{2}-512) \mathrm{r}^{4}-480 \mathrm{Rr}^{3}$.

## Proof.

1) From Theorem 11 it results that:

$$
\begin{aligned}
& \sum a^{3} \leq\left(r+\sqrt{4 R^{2}+r^{2}}\right)\left(8 R^{2}-2 r^{2}-2 r \sqrt{4 R^{2}+r^{2}}\right) \\
&= 8 R^{2} r-2 r^{3}-2 r^{2} \sqrt{4 R^{2}+r^{2}}+8 R^{2} \sqrt{4 R^{2}+r^{2}} \\
&-2 r^{2} \sqrt{4 R^{2}+r^{2}}-8 R^{2} r-2 r^{3} \\
&=\left(8 R^{2}-4 r^{2}\right) \sqrt{4 R^{2}+r^{2}}-4 r^{3} \\
& \leq\left(8 R^{2}-4 r^{2}\right)[2 R+(3-2 \sqrt{2}) r]-4 r^{3} \\
&=16 r^{3}+(24-16 \sqrt{2}) R^{2} r-8 R r^{2}-(12-8 \sqrt{2}) r^{3}-4 r^{3}
\end{aligned}
$$

which is equivalent with inequality from the statement.
2) From Theorem 11 it results that

$$
\sum a^{3} b \leq\left(r+\sqrt{4 R^{2}+r^{2}}\right)^{2}\left(8 R^{2}-4 r^{2}\right)
$$

and

$$
\left(r+\sqrt{4 R^{2}+r^{2}}\right)^{2} \leq 4 R^{2}+4 R r+(8-4 \sqrt{2}) r^{2}
$$

It results that:

$$
\begin{aligned}
\sum a^{3} b \leq & {\left[4 R^{2}+4 R r+(8-4 \sqrt{2}) r^{2}\right]\left(8 R^{2}-4 r^{2}\right) } \\
= & 32 R^{4}-16 R^{2} r^{2}+32 R^{3} r-16 R r^{3}+(64-32 \sqrt{2}) R^{2} r^{2} \\
& -(32-16 \sqrt{2}) r^{4}
\end{aligned}
$$

which is equivalent with the inequality from the statement.
3) We prove that:

$$
\begin{aligned}
\sum a^{3} b & \geq 352 R^{2} r^{2}+208 r^{4}-240 r^{3} \sqrt{4 R^{2}+r^{2}} \\
& \geq 352 R^{2} r^{2}+208 r^{4}-240 r^{3}[2 R+(3-2 \sqrt{2}) r] \\
& =352 R^{2} r^{2}+208 r^{4}-480 R r^{3}-(720-480 \sqrt{2}) r^{4}
\end{aligned}
$$

which is equivalent with the inequality from the statement.

## References

[1] W. J. Blundon, R. H. Eddy, Problem 488, Nieuw Arch. Wiskunde, 26 (1978).
[2] L. Fejes-Tóth, Inequalities concerning poligons and polyedra, Duke Math. J., 15 (1948), 817-822.
[3] T. Ovidiu, N. Pop Minculete, M. Bencze, An introduction to quadrilateral geometry, Editura Didactică şi Pedagogică, Bucureşti, 2013.
[4] Octogon Mathematical Magazine (1993-2013)

Received: 25 September 2013

# Approximation by limit q-Bernstein operator 

Zoltán Finta<br>Babeş-Bolyai University<br>Department of Mathematics<br>1, M. Kogălniceanu st., 400084<br>Cluj-Napoca, Romania<br>email: fzoltan@math.ubbcluj.ro

Dedicated to the memory of Professor Antal Bege


#### Abstract

We establish quantitative estimates for the limit q-Bernstein operator introduced in [3], via the second order Ditzian-Totik modulus of smoothness.


## 1 Introduction

The q-Bernstein operators were introduced by Phillips in [8] and they generalize the well-known Bernstein operators. A survey of the obtained results and references concerning q-Bernstein operators can be found in [6]. It is worth mentioning that the first generalization of the Bernstein operators based on q-integers was obtained by Lupaş [4].

Let $q>0$. For each nonnegative integer $k$, the $q$-integers $[k] \equiv[k]_{q}$ and the $q$-factorials [k]! are defined by

$$
[k]=\left\{\begin{array}{rll}
1+q+\cdots+q^{k-1}, & \text { if } k \geq 1 \\
0, & \text { if } k=0
\end{array}\right.
$$

2010 Mathematics Subject Classification: 41A25, 41A36
Key words and phrases: q-Bernstein operators, limit q-Bernstein operator, Ditzian-Totik modulus of smoothness, K-functional
and

$$
[k]!=\left\{\begin{array}{rl}
{[1][2] \ldots[k],} & \text { if } \quad k \geq 1 \\
1, & \text { if }
\end{array} \quad k=0\right.
$$

For integers $0 \leq k \leq n$, the $q$-binomial coefficients are defined by

$$
\left[\begin{array}{l}
n \\
k
\end{array}\right]=\frac{[n]!}{[k]![n-k]!}
$$

The $q$-Bernstein operators $B_{n, q}: C[0,1] \rightarrow C[0,1]$ are given by

$$
\begin{equation*}
\left(B_{n, q} f\right)(x) \equiv B_{n, q}(f, x)=\sum_{k=0}^{n} f\left(\frac{[k]}{[n]}\right) p_{n, k}(q, x) \tag{1}
\end{equation*}
$$

where $\mathrm{n}=1,2, \ldots, 0<\mathrm{q} \leq 1, x \in[0,1]$ and

$$
p_{n, k}(q, x)=\left[\begin{array}{l}
n \\
k
\end{array}\right] x^{k}(1-x)(1-x q) \ldots\left(1-x q^{n-k-1}\right)
$$

for $k=0,1, \ldots, n$ (an empty product is taken to be equal 1 ). For $q=1$ we recover the Bernstein operators. In [8], it is proved the uniform convergence of $B_{n, q_{n}} f$ to $f$ on $[0,1]$, as $n \rightarrow \infty$, when $q=q_{n} \in(0,1)$ and $q_{n} \rightarrow 1$ as $n \rightarrow \infty$.

Let $q \in(0,1)$ and $f \in C[0,1]$ be given. Il'inskii and Ostrovska proved in [3] that the sequence $\left\{B_{n, q}(f, x)\right\}$ converges to $B_{\infty, q}(f, x)$ as $n \rightarrow \infty$, uniformly for $x \in[0,1]$, where the limit $q$-Bernstein operator $\mathrm{B}_{\infty, \mathrm{q}}: \mathrm{C}[0,1] \rightarrow \mathrm{C}[0,1]$ is defined by

$$
\begin{align*}
& \left(B_{\infty, q} f\right)(x) \equiv B_{\infty, q}(f, x) \\
& \quad=\left\{\begin{aligned}
\sum_{k=0}^{\infty} f\left(1-q^{k}\right) \frac{x^{k}}{(1-q)^{k}[k]!} \prod_{s=0}^{\infty}\left(1-x q^{s}\right), & \text { if } 0 \leq x<1 \\
f(1), & \text { if } x=1 .
\end{aligned}\right. \tag{2}
\end{align*}
$$

The approximation of continuous functions $f$ by $B_{\infty, q} f$ as $q \nearrow 1$, has been investigated by Videnskii in [9]. We cite the following result of Videnskii. If $0<\mathrm{q}<1, \mathrm{x} \in[0,1]$ and $\mathrm{f} \in \mathrm{C}[0,1]$, then

$$
\begin{equation*}
\left|B_{\infty, q}(f, x)-f(x)\right| \leq 2 \omega\left(f, \frac{1}{2} \sqrt{1-q}\right) \tag{3}
\end{equation*}
$$

where $\omega(f, \delta)=\sup \{|f(x)-f(y)|: x, y \in[0,1],|x-y| \leq \delta\}$ is the usual modulus of continuity of $f$. For the second modulus of smoothness of $f$, defined by

$$
\omega^{2}(f, \delta)=\sup _{0<h \leq \delta} \sup _{x \in[0,1-2 h]}|f(x+2 h)-2 f(x+h)+f(x)|
$$

Wang obtained the following estimate (see [10] and [11]):

$$
\begin{equation*}
\left|B_{n, q}(f, x)-B_{\infty, q}(f, x)\right| \leq C \omega^{2}\left(f, \sqrt{q^{n}}\right) \tag{4}
\end{equation*}
$$

where $n=1,2, \ldots, x \in[0,1], 0<q<1$ and $f \in C[0,1]$. Here we mention that $C>0$ is a constant independent of $n, x$ and $q$, which can be different at each occurrence.

The goal of the paper is to establish quantitative results for the rate of convergence of (2), obtaining similar estimates to (3) and (4). In our estimates we shall use the second order Ditzian-Totik modulus of smoothness of $f$, defined by

$$
\omega_{\varphi}^{2}(f, \delta)=\sup _{0<h \leq \delta} \sup _{x \pm h \varphi(x) \in[0,1]}|f(x+h \varphi(x))-2 f(x)+f(x-h \varphi(x))|
$$

where $\varphi(x)=\sqrt{x(1-x)}, x \in[0,1]$ (for details see [1]). Further, we consider the following K-functional:

$$
K_{2, \varphi}(f, \delta)=\inf _{g \in W^{2}(\varphi)}\left\{\|f-g\|+\delta\left\|\varphi^{2} g^{\prime \prime}\right\|\right\}
$$

where $\|\cdot\|$ denotes the uniform norm on $C[0,1]$ and $W^{2}(\varphi)=\left\{g \in C[0,1]: g^{\prime} \in\right.$ $\left.A C_{l o c}[0,1], \varphi^{2} g^{\prime \prime} \in \mathrm{C}[0,1]\right\} ; \mathrm{g}^{\prime} \in A C_{\text {loc }}[0,1]$ means that $g$ is differentiable such that $g^{\prime}$ is absolutely continuous on every interval $[a, b] \subset[0,1]$. It is known (see $[1,(2.1 .4)])$ that $\omega_{\varphi}^{2}(f, \sqrt{\delta})$ and $K_{2, \varphi}(f, \delta)$ are equivalent, i.e. there exists $C>0$ such that

$$
\begin{equation*}
\mathrm{C}^{-1} \omega_{\varphi}^{2}(\mathrm{f}, \sqrt{\delta}) \leq \mathrm{K}_{2, \varphi}(\mathrm{f}, \delta) \leq \mathrm{C} \omega_{\varphi}^{2}(\mathrm{f}, \sqrt{\delta}) \tag{5}
\end{equation*}
$$

## 2 Main results

Theorem 1 There exists $C>0$ such that

$$
\left\|B_{\infty, q} f-f\right\| \leq C \omega_{\varphi}^{2}(f, \sqrt{1-q})
$$

for all $\mathrm{f} \in \mathrm{C}[0,1]$ and $\mathrm{q} \in(0,1)$. Consequently, $\mathrm{B}_{\infty, \mathrm{q}} \mathrm{f}$ converges uniformly to f on $[0,1]$ as $\mathrm{q} \nearrow 1$.

Proof. By $[9,(7.7)-(7.8)]$, we have $B_{\infty, q}(1, x)=1$ and $B_{\infty, q}(t, x)=x$. For $g \in W^{2}(\varphi)$, by Taylor's formula:

$$
g(t)=g(x)+g^{\prime}(x)(t-x)+\int_{x}^{t}(t-u) g^{\prime \prime}(u) d u, \quad t, x \in[0,1]
$$

we get

$$
B_{\infty, q}(g, x)-g(x)=B_{\infty, q}\left(\int_{x}^{t}(t-u) g^{\prime \prime}(u) d u, x\right)
$$

Using the inequality

$$
\begin{equation*}
\left|\int_{x}^{t}(t-u) g^{\prime \prime}(u) d u\right| \leq(t-x)^{2} \varphi^{-2}(x)\left\|\varphi^{2} g^{\prime \prime}\right\| \tag{6}
\end{equation*}
$$

(see $\left[1\right.$, Lemma 9.6.1]) and $B_{\infty, q}\left((t-x)^{2}, x\right)=(1-q) \varphi^{2}(x)($ see $[9,(7.12)])$, we find

$$
\begin{align*}
\left|B_{\infty, q}(g, x)-g(x)\right| & \leq B_{\infty, q}\left(\left|\int_{x}^{t}\right| t-u \| g^{\prime \prime}(u)|d u|, x\right) \\
& \leq B_{\infty, q}\left((t-x)^{2}, x\right) \varphi^{-2}(x)\left\|\varphi^{2} g^{\prime \prime}\right\| \\
& =(1-q)\left\|\varphi^{2} g^{\prime \prime}\right\| \tag{7}
\end{align*}
$$

On the other hand, by $(2)$ and $B_{\infty, q}(1, x)=1$, we obtain $\left|B_{\infty, q}(f, x)\right| \leq$ $\|f\| B_{\infty, q}(1, x)=\|f\|$, i.e.

$$
\begin{equation*}
\left\|\mathrm{B}_{\infty, \mathrm{q}} \mathrm{f}\right\| \leq\|\mathrm{f}\| \tag{8}
\end{equation*}
$$

for all $f \in C[0,1]$. Now, in view of (7) and (8), we get

$$
\begin{aligned}
\left\|B_{\infty, q} f-f\right\| & \leq\left\|B_{\infty, q} f-B_{\infty, q} g\right\|+\left\|B_{\infty, q} g-g\right\|+\|g-f\| \\
& \leq 2\|f-g\|+(1-q)\left\|\varphi^{2} g^{\prime \prime}\right\| \\
& \leq 2\left\{\|f-g\|+(1-q)\left\|\varphi^{2} g^{\prime \prime}\right\|\right\}
\end{aligned}
$$

Taking the infimum on the right-hand side over all $\mathrm{g} \in \mathrm{W}^{2}(\varphi)$ and using (5), we get the assertion of our theorem.

Remark 1 The main result of [2] provides an estimate for positive linear operators that preserve linear functions. The result was improved in $[7,(2.138)]$, which implies for the limit $q$-Bernstein operator that

$$
\left\|B_{\infty, q} f-f\right\| \leq \frac{5}{2} \omega_{\varphi}^{2}(f, \sqrt{1-q}), \quad \text { where } \quad \frac{3}{4} \leq q<1
$$

Theorem 2 Let $\mathrm{q} \in(0,1)$ be given. Then there exists $\mathrm{C}>0$ such that

$$
\left\|B_{n, q} f-B_{\infty, q} f\right\| \leq \frac{C}{q(1-q)} \omega_{\varphi}^{2}\left(f, \sqrt{q^{n}}\right)
$$

for all $\mathrm{f} \in \mathrm{C}[0,1]$ and $\mathrm{n}=1,2, \ldots$.
Proof. Let $g \in W^{2}(\varphi)$ and $x \in[0,1]$. Then, by [5, (3.2)], we have

$$
\begin{equation*}
B_{n, q}(g, x)-B_{n+1, q}(g, x)=\sum_{k=1}^{n} a_{n, k}(g) p_{n+1, k}(q, x) \tag{9}
\end{equation*}
$$

where

$$
\begin{align*}
a_{n, k}(g)= & \frac{[n+1-k]}{[n+1]} g\left(\frac{[k]}{[n]}\right)+q^{n+1-k} \frac{[k]}{[n+1]} g\left(\frac{[k-1]}{[n]}\right) \\
& -g\left(\frac{[k]}{[n+1]}\right) . \tag{10}
\end{align*}
$$

By Taylor's formula, we find

$$
\begin{aligned}
g\left(\frac{[k]}{[n]}\right)= & g\left(\frac{[k]}{[n+1]}\right)+\left(\frac{[k]}{[n]}-\frac{[k]}{[n+1]}\right) g^{\prime}\left(\frac{[k]}{[n+1]}\right) \\
& +\int_{[k] /[n+1]}^{[k] /[n]}\left(\frac{[k]}{[n]}-u\right) g^{\prime \prime}(u) d u
\end{aligned}
$$

and

$$
\begin{aligned}
g\left(\frac{[k-1]}{[n]}\right)= & g\left(\frac{[k]}{[n+1]}\right)+\left(\frac{[k-1]}{[n]}-\frac{[k]}{[n+1]}\right) g^{\prime}\left(\frac{[k]}{[n+1]}\right) \\
& +\int_{[k] /[n+1]}^{[k-1] /[n]}\left(\frac{[k-1]}{[n]}-u\right) g^{\prime \prime}(u) d u,
\end{aligned}
$$

respectively. Hence, by (10),

$$
\begin{aligned}
a_{n, k}(g)= & \frac{[n+1-k]}{[n+1]} g\left(\frac{[k]}{[n]}\right)+q^{n+1-k} \frac{[k]}{[n+1]} g\left(\frac{[k-1]}{[n]}\right) \\
& -\frac{[n+1-k]+q^{n+1-k}[k]}{[n+1]} g\left(\frac{[k]}{[n+1]}\right) \\
= & \frac{[n+1-k]}{[n+1]}\left(\frac{[k]}{[n]}-\frac{[k]}{[n+1]}\right) g^{\prime}\left(\frac{[k]}{[n+1]}\right)
\end{aligned}
$$

$$
\begin{align*}
& +\frac{[n+1-k]}{[n+1]} \int_{[k] /[n+1]}^{[k] /[n]}\left(\frac{[k]}{[n]}-u\right) g^{\prime \prime}(u) d u \\
& +\frac{q^{n+1-k}[k]}{[n+1]}\left(\frac{[k-1]}{[n]}-\frac{[k]}{[n+1]}\right) g^{\prime}\left(\frac{[k]}{[n+1]}\right) \\
& +\frac{q^{n+1-k}[k]}{[n+1]} \int_{[k] /[n+1]}^{[k-1] /[n]}\left(\frac{[k-1]}{[n]}-u\right) g^{\prime \prime}(u) d u \\
& =\frac{[n+1-k]}{[n+1]} \int_{[k] /[n+1]}^{[k] /[n]}\left(\frac{[k]}{[n]}-u\right) g^{\prime \prime}(u) d u \\
& +\frac{q^{n+1-k}[k]}{[n+1]} \int_{[k] /[n+1]}^{[k-1] /[n]}\left(\frac{[k-1]}{[n]}-u\right) g^{\prime \prime}(u) d u, \tag{11}
\end{align*}
$$

because

$$
\begin{aligned}
& \frac{[n+1}{[n+1]}\left(\frac{k]}{[n]}-\frac{[k]}{[n+1]}\right)+\frac{q^{n+1-k}[k]}{[n+1]}\left(\frac{[k-1]}{[n]}-\frac{[k]}{[n+1]}\right) \\
& \quad= \frac{[k]}{[n][n+1]^{2}}\{[n+1-k]([n+1]-[n]) \\
&\left.+q^{n+1-k}([k-1][n+1]-[k][n])\right\} \\
&=\frac{[k]}{[n][n+1]^{2}}\left\{[n+1-k] q^{n}+q^{n+1-k}\left(-q^{k-1}[n+1-k]\right)\right\} \\
& \quad= 0 .
\end{aligned}
$$

In view of (6) and (11), we have

$$
\begin{aligned}
\left|a_{n, k}(g)\right| \leq & \frac{[n+1-k]}{[n+1]}\left(\frac{[k]}{[n]}-\frac{[k]}{[n+1]}\right)^{2} \varphi^{-2}\left(\frac{[k]}{[n+1]}\right)\left\|\varphi^{2} g^{\prime \prime}\right\| \\
& +\frac{q^{n+1-k}[k]}{[n+1]}\left(\frac{[k-1]}{[n]}-\frac{[k]}{[n+1]}\right)^{2} \varphi^{-2}\left(\frac{[k]}{[n+1]}\right)\left\|\varphi^{2} g^{\prime \prime}\right\| \\
= & \left\{\frac{[n+1-k][k]([n+1]-[n])^{2}}{[n]^{2}[n+1]([n+1]-[k])}\right. \\
& \left.+\frac{q^{n+1-k}([k-1][n+1]-[k][n])^{2}}{[n]^{2}[n+1]([n+1]-[k])}\right\}\left\|\varphi^{2} g^{\prime \prime}\right\| \\
= & \left\{\frac{[n+1-k][k] q^{2 n}}{[n]^{2}[n+1] q^{k}[n+1-k]}\right.
\end{aligned}
$$

$$
\begin{aligned}
& \left.+\frac{q^{n+1-k}\left(-q^{k-1}[n+1-k]\right)^{2}}{[n]^{2}[n+1] q^{k}[n+1-k]}\right\}\left\|\varphi^{2} g^{\prime \prime}\right\| \\
= & \frac{q^{n-1}}{[n]^{2}[n+1]}\left\{q^{n+1-k}[k]+[n+1-k]\right\}\left\|\varphi^{2} g^{\prime \prime}\right\| \\
= & \frac{q^{n-1}}{[n]^{2}}\left\|\varphi^{2} g^{\prime \prime}\right\| \leq q^{n-1}\left\|\varphi^{2} g^{\prime \prime}\right\| .
\end{aligned}
$$

Hence, by (9) and $B_{n+1, q}(1, x)=1$ (see $[9,(2.5)]$ ), we find

$$
\left|B_{n, q}(g, x)-B_{n+1, q}(g, x)\right| \leq q^{n-1}\left\|\varphi^{2} g^{\prime \prime}\right\|
$$

for all $x \in[0,1]$. This implies that

$$
\begin{align*}
\| B_{n, q} g- & B_{n+p, q} g \| \\
\leq & \left\|B_{n, q} g-B_{n+1, q} g\right\|+\left\|B_{n+1, q} g-B_{n+2, q} g\right\| \\
& +\cdots+\left\|B_{n+p-1, q} g-B_{n+p, q} g\right\| \\
\leq & \left(q^{n-1}+q^{n}+\cdots+q^{n+p-2}\right)\left\|\varphi^{2} g^{\prime \prime}\right\| \\
\leq & \frac{q^{n-1}}{1-q}\left\|\varphi^{2} g^{\prime \prime}\right\| \tag{12}
\end{align*}
$$

for $n, p=1,2, \ldots$ In conclusion $\left\{B_{n, q} g\right\}$ is a Cauchy-sequence in $C[0,1]$, so $\left\{B_{n, q} g\right\}$ converges to $B_{\infty, q} g$ as $n \rightarrow \infty$ (see also [3]). Now let $p \rightarrow \infty$ in (12). Then we obtain

$$
\begin{equation*}
\left\|B_{n, q} g-B_{\infty, q} g\right\| \leq \frac{q^{n}}{q(1-q)}\left\|\varphi^{2} g^{\prime \prime}\right\| \tag{13}
\end{equation*}
$$

Further, by (1) and $B_{n, q}(1, x)=1$ (see $[9,(2.5)]$ ), we obtain $\left|B_{n, q}(f, x)\right| \leq$ $\|f\| B_{n, q}(1, x)=\|f\|$, i.e.

$$
\begin{equation*}
\left\|B_{n, q} f\right\| \leq\|f\| \tag{14}
\end{equation*}
$$

for all $f \in C[0,1]$. Then (14), (8) and (13) imply that

$$
\begin{aligned}
\left\|B_{n, q} f-B_{\infty, q} f\right\| \leq & \left\|B_{n, q} f-B_{n, q} g\right\|+\left\|B_{n, q} g-B_{\infty, q} g\right\| \\
& +\left\|B_{\infty, q} g-B_{\infty, q} f\right\| \\
\leq & 2\|f-g\|+\frac{q^{n}}{q(1-q)}\left\|\varphi^{2} g^{\prime \prime}\right\| \\
\leq & \frac{2}{q(1-q)}\left\{\|f-g\|+q^{n}\left\|\varphi^{2} g^{\prime \prime}\right\|\right\}
\end{aligned}
$$

Taking the infimum on the right-hand side over all $g \in W^{2}(\varphi)$ and using (5), we get the assertion of our theorem.

Remark 2 Because $\omega_{\varphi}^{2}(f, \delta) \leq \mathrm{C} \omega^{2}(f, \delta) \leq 2 \mathrm{C} \omega(f, \delta)$ (for details see [1]), we obtain, in view of Theorem 1 and Theorem 2, the following weaker estimates:

$$
\left|B_{\infty, q}(f, x)-f(x)\right| \leq C \omega(f, \sqrt{1-q})
$$

and

$$
\left|B_{\infty, q}(f, x)-B_{n, q}(f, x)\right| \leq \frac{C}{q(1-q)} \omega^{2}\left(f, \sqrt{q^{n}}\right)
$$

## References

[1] Z. Ditzian, V. Totik, Moduli of smoothness, Springer, Berlin, 1987.
[2] I. Gavrea, H. Gonska, R. Păltănea, G. Tachev, General estimates for the Ditzian-Totik modulus, East J. Approx., 9 (2) (2003), 175-194.
[3] A. Il'inskii, S. Ostrovska, Convergence of generalized Bernstein polynomials, J. Approx. Theory, 116 (2002), 100-112.
[4] A. Lupaş, A q-analogue of the Bernstein operator, Babeş-Bolyai University, Seminar on Numerical and Statistical Calculus, 9 (1987), 85-92.
[5] H. Oruç, G. M. Phillips, A generalization of the Bernstein polynomials, Proc. Edinb. Math. Soc., 42 (1999), 403-413.
[6] S. Ostrovska, The first decade of the q-Bernstein polynomials: results and perspectives, J. Math. Anal. Approx. Theory, 2 (1) (2007), 35-51.
[7] R. Păltănea, Approximation theory using positive linear operators, Birkhäuser, Boston, 2004.
[8] G. M. Phillips, Bernstein polynomials based on the q-integers, Ann. Numer. Math., 4 (1997), 511-518.
[9] V. S. Videnskii, On some classes of q-parametric positive operators, $O p$ erator Theory: Advances and Applications, 158 (2005), 213-222.
[10] H. Wang, Korovkin-type theorem and applications, J. Approx. Theory, 132 (2005), 258-264.
[11] H. Wang, F. Meng, The rate of convergence of q-Bernstein polynomials for $0<\mathrm{q}<1$, J. Approx. Theory, 136 (2005), 151-158.

# Galois covering and smash product of skew categories 

Emil Horobeţ<br>Technical University of Eindhoven<br>email: e.horobet@tue.nl<br>Dedicated to the memory of Professor Antal Bege


#### Abstract

In this paper we give a new proof of the famous result of E. L. Green [3], that gradings of a finite, path connected quiver are in one-to-one correspondence with Galois coverings. Namely we prove that the inverse construction to the skew group construction has as many solutions as the number of different gradings on the starting quiver.


## 1 Introduction

In mathematics, a k-category or an Abelian category is a category in which morphisms and objects can be added. The motivation for k-categories originated from examining the category of Abelian groups, $A b$. The theory rises from a tentative attempt to unify several cohomology theories by A. Grothendieck.

In this paper we examine skew categories and their connection to skew group algebras. The paper has two parts. In the first part we recall the basic notions and results of this topic, for this we use as a basis literature [1] and [2]. We present the categorical machinery developed by the above mentioned authors among some reformulations of several coherence results, to fit better in our context. In the second part we give a new proof of the famous result of E . L. Green [3], that gradings of a finite, path connected quiver are in one-toone correspondence with Galois coverings. More precisely we prove that the inverse construction to the skew group construction has as many solutions as the number of different gradings on the starting quiver.

[^1]Key words and phrases: smash product, Galois covering, skew group algebras

### 1.1 Basic notions

From now on we deal with small categories over a commutative field $k$, this means that the objects $C_{0}$ form a set and not a class, and the morphisms consist of modules over $k$. We set $G$ to be an arbitrary group. We define the notions of group action and grading on k-categories by [2].

Definition 1 A G-category is a category C with the propriety that G acts on the set of objects, in other words the elements of G are k -module morphisms, such that the following hold: for all $s \in G$, and for all $x, y \in C_{0}$ for $s:{ }_{y} C_{x} \rightarrow_{s y}$ $\mathrm{C}_{\mathrm{sx}}$, we have that

- $\mathbf{s}(\mathrm{gf})=(\mathrm{sg})(\mathrm{sf})$, if f and g can be composed in C ;
- If $\mathrm{t}, \mathrm{s} \in \mathrm{G}$ and f is a morphism in C , then $(\mathrm{ts}) \mathrm{f}=\mathrm{t}(\mathrm{sf})$;
- $1 \mathrm{f}=\mathrm{f}, 1 \in \mathrm{G}$ the identity element.

In other words $G$ is a group of autofunctors of $C$.
Since we defined a G-action on our categories it makes sense to talk about graded categories. These gradings will play a crucial role in the inverse construction we deal with in the second section.

Definition 2 A G-graded category is a category C for which the following hold:

- For all $x, y \in C_{0}$ we have ${ }_{y} \mathrm{C}_{\mathrm{x}}=\bigoplus_{\mathrm{s} \in \mathrm{G}}\left(\mathrm{y}_{\mathrm{y}} \mathrm{C}_{\mathrm{x}}^{\mathrm{s}}\right)$ and

$$
{ }_{z} C_{y}^{t} y C_{x}^{s} C_{y} C_{x}^{t s}
$$

- ${ }_{x} 1_{x} \in_{x} C x^{1}$.

For the definition of Galois covering of categories we must define first what a quotient category is. For this we cite [2], Definition 2.1.

Definition 3 If C is a free G -category over k , then the objects of the quotient category $\mathrm{C} / \mathrm{G}$ are the G -orbits of $\mathrm{C}_{0}$ and if $\alpha, \beta$ are two G -orbits, then the morphisms between them are:

$$
\beta(\mathrm{C} / \mathrm{G})_{\alpha}=\left(\bigoplus_{x \in \alpha, y \in \beta}\left({ }_{y} C_{x}\right)\right) / \mathrm{G} .
$$

Now let $\mathrm{p}: \mathrm{C} \rightarrow \mathrm{C} / \mathrm{G}$ be the projection functor, then we call p the Galois covering of C with Galois group G .

Similarly to the above construction we are interested in the categorical definition of skew group algebras, namely skew categories ([2], Definition 2.3).

Definition 4 Let C be a G-category. Then the objects of the skew category $\mathrm{C}[\mathrm{G}]$ are the objects of the category C , so we have $(\mathrm{C}[\mathrm{G}])_{0}=\mathrm{C}_{0}$ and the morphisms between them are

$$
y(C[G])_{x}=\bigoplus_{s \in G}\left({ }_{y} C_{s x}\right)
$$

It is natural to ask if the above construction gives back the classical notion of skew group algebras. For this we cite a coherence result with the usual skew algebra construction ([2], Proposition 2.4).

Proposition 1 Let G be a finite group and let C be G-category over k, with finite number of objects. Let $\mathbf{a}(\mathrm{C})$ be the k -algebra associated to C , namely

$$
\mathrm{a}(\mathrm{C})=\bigoplus_{x, y \in \mathrm{C}_{0}} y \mathrm{C}_{x},
$$

provided with the matrix product induced by the composition of morphisms. Then we have that

$$
a(C[G]) \cong a(C)[G] .
$$

Now let us see how Galois coverings and skew group algebras are related. The following theorem is very important for the theory ([2], Theorem 2.8) and from now on we will use it implicitly without referring to it.

Proposition 2 Let C be a free G-category over k . The quotient category C/G and the skew category $\mathrm{C}[\mathrm{G}]$ are equivalent.

The main goal of this part is to present the necessary tools for developing the inverse construction to taking the quotient of a category, this will be the smash product category ([2], Definition 3.1).

Definition 5 Let G be a group and let C be a G-graded category over k . Then the smash product category C\#G has object set $\mathrm{C}_{0} \times \mathrm{G}$. Let $(\mathrm{x}, \mathrm{s}),(\mathrm{y}, \mathrm{t}) \in$ $\mathrm{C}_{0} \times \mathrm{G}$ be two objects. The k -module morphisms are defined as follows:

$$
{ }_{(y, t)}(C \# G)_{(x, s)}=y_{y} C_{x}^{t^{-1} s}
$$

It is natural to ask again if this construction gives back the classical notion of the smash product of algebras. For this we present a coherence result in a form to serve better our further goals.

Proposition 3 Let G be a finite group and let C be a G-category, then the k -algebras $\mathrm{a}(\mathrm{C}) \# \mathrm{G}$ and $\mathrm{a}(\mathrm{C} \# \mathrm{G})$ are Morita equivalent.

Since this is not the classical statement regarding smash product categories we present a proof for it.

Proof. If C is a G-category, then every morphism space ${ }_{y} C_{x}$ is a G-module, but $G$ being finite one can also regard these spaces as $(\mathrm{kG})^{*}$ modules. In this setting $C$ can be thought as a (kG)*-module category.

Now by Theorem 2.9 of [1] we have that the k-categories C\#G and C\#(kG)* are Morita equivalent. Moreover we can derive from this that $a(C \# G)$ and $\mathrm{a}\left(\mathrm{C} \#(\mathrm{kG})^{*}\right)$ are Morita equivalent as k -algebras.

Combining this with Proposition 2.3 from [1], which claims that the kalgebras $\mathfrak{a}(\mathrm{C}) \#(\mathrm{kG})^{*}$ and $\mathfrak{a}\left(\mathrm{C} \#(\mathrm{kG})^{*}\right)$ are isomorphic, we get that $\mathfrak{a}(\mathrm{C}) \# \mathrm{G}$ and $\mathrm{a}(\mathrm{C} \# \mathrm{G})$ are Morita equivalent.

A last definition before we reach to the main duality theorems of this section, is of a matrix category ([1], Definition 4.1).

Definition 6 Let C be ak-category and let n be a sequence of positive integers $\left(\mathrm{n}_{\mathrm{x}}\right)_{\mathrm{x} \in \mathrm{C}_{0}}$. The object set of the matrix category $\mathrm{M}_{\mathrm{n}}(\mathrm{C})$ remains the same objects of C . The set of morphisms from x to y is the vector space of $\mathrm{n}_{\mathrm{x}}$ columns and $n_{y}$ rows rectangular matrices with entries in ${ }_{y} C_{x}$. Composition of morphisms is given by the matrix product combined with the composition in C .

A classical way of relating the matrix categories to the corresponding matrix algebras is to consider single object categories provided by an algebra $A$ and then proving that the matrix category has one object with endomorphism algebra precisely the usual algebra of matrices $M_{n}(A)$. Unfortunately this approach is not sufficient for our further goal, so we need to develop a different correspondence between these categorical and ring theoretical objects. For this we have the following lemma.

Lemma 1 Let C be a k -category and let n be a positive integer, then we have the following k -algebra isomorphism

$$
a\left(M_{n}(C)\right) \cong M_{n}(k C)
$$

where in the right hand side kC is regarded as the path algebra of the underlying quiver of $C$.

Proof. Let us examine carefully the construction of the morphism spaces of the matrix category, we have that

$$
\begin{aligned}
a\left(M_{n}(C)\right) & =\bigoplus_{x, y \in M_{n}(C)_{0}} y\left(M_{n}(C)\right)_{x}=\bigoplus_{x, y \in C_{0}} M_{n}\left({ }_{y} C_{x}\right)= \\
& =M_{n}\left(\bigoplus_{x, y \in C_{0}} y C_{x}\right)=M_{n}(k C)
\end{aligned}
$$

Here we consider the vertices as the identity morphisms on the corresponding object, hence the set of vertices is a subset of the set of all morphisms. In this respect we can consider $\mathfrak{a}(\mathrm{C})$ isomorphic to the path algebra kC .
Going back for a moment to the matrix categories we want to recall the following equivalence ([1], Corollary 4.5).

Proposition 4 Let C be a k-category and n a positive integer, then C and $\mathrm{M}_{\mathrm{n}}(\mathrm{C})$ are Morita equivalent.

Now the last statement of this section is the categorical version of the CohenMongomery duality ([2], Proposition 3.2).

Theorem 1 Let C be a G-graded category over k . Then the category (C\#G)[G] is equivalent to C .

## 2 The inverse construction

Now that we presented the categorical machinery developed for skew categories and smash products, we pass to the main theorem of this paper, namely the inverse construction to the skew group construction. From now on we consider finite, path connected quivers as categories over $k$ : the objects are the vertices of the quiver and morphisms between two vertices are free $k$-modules having a basis given by the paths between these vertices.

Theorem 2 Let C be a finite, path connected quiver, and let G be a group acting on it. Given a G-grading on C, we have that the skew group algebra $\left(\mathrm{kC}_{\mathrm{G}}\right)[\mathrm{G}]$ and the path algebra kC are Morita equivalent, where $\mathrm{C}_{\mathrm{G}}$ is the quiver corresponding to $\mathrm{C} \# \mathrm{G}$.

Proof. We are considering C as a k-category, then by the Cohen-Mongomery duality (Theorem 1) we have the following equivalence of categories

$$
(\mathrm{C} \# \mathrm{G})[\mathrm{G}] \cong \mathrm{C} .
$$

Translating this to the language of k-algebras, via the functor $a$, we get that

$$
a((C \# G)[G]) \cong a(C)
$$

as k -algebras. Now by the remark in the proof of Lemma 1 we can consider $a(C)$ to be the path algebra $k C$.

Applying the coherence property of the skew group construction (Proposition 1), we get the following isomorphism of algebras

$$
\mathrm{a}(\mathrm{C} \# \mathrm{G})[\mathrm{G}] \cong \mathrm{kC}
$$

From this point, by applying the coherence result of the smash product (Proposition 3), we pass to Morita equivalences. So we get that ( $\mathrm{a}(\mathrm{C}) \# \mathrm{G})[\mathrm{G}]$ is Morita equivalent to $k C$, where $a(C) \# G$ is a smash product of algebras.

Finally applying again the remark from Lemma 1, we get that $(\mathrm{k}(\mathrm{C} \# \mathrm{G}))_{[G]}$ is Morita equivalent to $k C$, here $k(C \# G)$ is viewed as the path algebra of the quiver corresponding to $\mathrm{C} \# \mathrm{G}$.

Now putting everything in our notation we get the expected result, that the skew group algebra $\left(\mathrm{kC}_{\mathrm{G}}\right)[\mathrm{G}]$ and the path algebra kC are Morita equivalent, where $\mathrm{C}_{\mathrm{G}}$ is the quiver corresponding to $\mathrm{C} \# \mathrm{G}$.
One can see from the above result that each different grading of $C$ will lead to a different solution $\mathrm{C} \# \mathrm{G}$ to the inverse construction problem.

Finally we give an example to illustrate our result.
Example 1 Let $\mathrm{G}=\langle\mathrm{g}\rangle$ be a cyclic group of order two and let C be the following quiver


We can consider C to be a G-graded quiver by setting degree 1 for the elements $\{\mathrm{e}, \mathrm{f}\}$ and degree g for the elements $\{\alpha, \beta\}$. In this case the quiver $\mathrm{C}_{\mathrm{G}}$, corresponding to the smash product of C and G is the following


So we get that the skew group algebra $\left(\mathrm{kC}_{\mathrm{G}}\right)[\mathrm{G}]$ is Morita equivalent to the path algebra of C .

## Acknowledgements

This work was supported by a grant of the Ministry of National Education, CNCS-UEFISCDI, project number PN-II-ID-PCE-2012-4-0100.

## References

[1] C. Cibils, A. Solotar, Galois coverings, Morita equvalence and Smash extensions of categories over a field, Documenta Mathematica, 11 (2006), 143-159.
[2] C. Cibils, E. N. Marcos, Skew Categories, Galois coverings and smash products of k-categories, Proc. Amer. Math. Soc., 134 (2005), 39-50.
[3] E. L. Green, Graphs with relations, coverings, and group-graded algebras, Trans. Amer. Math. Soc., 279 (1983), 297-310.

# Leader election in synchronous networks 

Antal Iványi<br>Faculty of Informatics<br>Eötvös Loránd University<br>Budapest, Hungary<br>email: tony@inf.elte.hu

Dedicated to the memory of my friend Professor Antal Bege


#### Abstract

Worst, best and average number of messages and running time of leader election algorithms of different distributed systems are analyzed. Among others the known characterizations of the expected number of messages for LCR algorithm and of the worst number of messages of Hirschberg-Sinclair algorithm are improved.


## 1 Introduction

We consider the problem of leader election in synchronous networks [11, 16, $30,43,59,92]$. The networks are modeled by directed graphs, the processors are called processes and are modeled as an automaton (see e.g. [11, 59]). In the case of the deterministic algorithms it is supposed that the processes have a unique identifier (UID).

The main topic of this paper is the presentation of leader election algorithms of different synchronous networks and their performance features.

It is known that if the processes are indistinguishable then there is no deterministic algorithm to solve the problem. For such anonymous or symmetric networks random algorithms are proposed by Itai and Rodeh [38, 39], by Ghaffni et al. [31], and by Kalpathi et al. [42].

Lower and upper bounds for the number of necessary messages or necessary bits are presented by Afek and Gafni, Attiya et al., Bodlaender, Frederickson and Lynch, Korach et al., and Loui et al. $[1,2,8,9,26,47,58]$.

2010 Mathematics Subject Classification: 68Q25, 68W10, 68W40
Key words and phrases: leader election, synchronous networks, analysis of algorithms, LCR, HS

The structure of the paper is as follows. After the introductory Section 1 in Section 2 the enumeration of some distributed systems is presented, then in Section 3 simple (as complete, chain, mesh and star networks), ring (unidirectional and bidirectional), special (such as De Bruijn, hypercube, Cayley, tree and recursively scalable networks) and general networks are analyzed.

## 2 Enumeration of labeled directed networks

Leader election requires that any process can inform any other process on its own data (e.g. on its own uid). In order to guarantee the participation of all processes we suppose that the investigated networks are strongly connected. It is worth to remark that there are also algorithms not requiring the strong connectedness, but these algorithms have also such output that the leader election is not solvable.

In this section we deal at first with the influence of the requirement of strong connectedness on the number of the tested networks, then with some simple networks such as complete network, star and chain.

### 2.1 Enumeration of connected and strongly connected networks

Let $D(n), C(n)$, and $S(n)$ denote the number of labeled simple, labeled simple weakly connected and labeled simple strongly connected digraphs, respectively.

The known simple formula

$$
\begin{equation*}
D(n)=2^{n(n-1)} \tag{1}
\end{equation*}
$$

gives $D(n)$. The values of $D(n)$, further $C(n) / D(n)$ and $S(n) / D(n)$ are shown in Table 1 for $n=1, \ldots, 15$. Values of $D(n)$ for $n=1, \ldots, 35$ can be found in [65].
In 2012 Critzer [17] proposed the following method to determine the number $\mathrm{C}(\mathrm{n})$ of the simple labeled weakly connected digraphs:

$$
\begin{equation*}
C(n)=D(n)-\frac{1}{n} \sum_{i=1}^{n-1} k\binom{n}{k} C(k) D(n-k) . \tag{2}
\end{equation*}
$$

Using (1) one can compute the $\mathrm{D}(\mathrm{n})$ values necessary to get the values of $C(n)$ from (2). E.g. (1) results $D(1)=1$ and then (2) gives $C(1)=1$. In a similar way $\mathrm{D}(2)=4$ and $\mathrm{C}(2)=3$, further $\mathrm{D}(3)=64$ and $\mathrm{C}(3)=54$.

| n | $\mathrm{D}(\mathrm{n})$ | $\mathrm{C}(\mathrm{n}) / \mathrm{D}(\mathrm{n})$ | $\mathrm{S}(\mathrm{n}) / \mathrm{D}(\mathrm{n})$ |
| ---: | ---: | ---: | ---: |
| 1 | 1 | 1.000000 | 1.00000 |
| 2 | 4 | 0.750000 | 0.25000 |
| 3 | 64 | 0.843750 | 0.28125 |
| 4 | 4096 | 0.936035 | 0.39209 |
| 5 | 1048576 | 0.979500 | 0.53890 |
| 6 | 4398046511104 | 0.998280 | 0.80106 |
| 7 | $\sim 441829$ | 0.994008 | 0.68431 |
| 8 | 72057594037927936 | 0.999511 | 0.88506 |
| 9 | $\sim 4.722366483 \cdot 10^{2}$ | 0.999863 | 0.93161 |
| 10 | $\sim 1.237940039 \cdot 10^{27}$ | 0.999962 | 0.96132 |
| 11 | $\sim 1.298074215 \cdot 10^{33}$ | 0.999990 | 0.97843 |
| 12 | $\sim 5.444517871 \cdot 10^{39}$ | 0.999997 | 0.98835 |
| 13 | $\sim 9.134385523 \cdot 10^{46}$ | 0.999999 | 0.99367 |
| 14 | $\sim 6.129982164 \cdot 10^{54}$ | 0.9999998 | 0.99659 |
| 15 | $\sim 1.645504557 \cdot 10^{60}$ | 0.99999994 | 0.99817 |

Table 1: Number $D(n)$ of simple labeled directed graphs and the ratios $C(n) /(D(n)$ and $S(n) / D(n)$.

Table 2 contains $C(n)$ for $n=1, \ldots, 15$. In $[66]$ the values for $n=16, \ldots, 35$ can be found.
V. A. Liskovets in 1969 [52, 100] proposed the following recursive formulas to compute $S(n)$ :

$$
\begin{align*}
& a(n)=n(n-1)-\sum_{i=1}^{n-1}\binom{n-1}{t-1} a(t)  \tag{3}\\
& \lambda_{t}(m)=2^{m(m+t-1)}-\sum_{k=0}^{m-1}\binom{m}{k} \lambda_{t}(k) \tag{4}
\end{align*}
$$

and

$$
\begin{equation*}
S(n)=a(n)+\sum_{i=1}\binom{n-1}{t-1} 2^{(m-1)(m-k)} \lambda_{t}(n-t) S(t) \tag{5}
\end{equation*}
$$

Using (3) and (4) one can compute the $a(n)$ and $\lambda(n)$ values necessary to get the values of $S(n)$ from (5).

| n | $\mathrm{C}(\mathrm{n})$ |
| :---: | :---: |
| 1 | 1 |
| 2 | 3 |
| 3 | 54 |
| 4 | 3834 |
| 5 | 1027080 |
| 6 | 1067308488 |
| 7 | 4390480193904 |
| 8 | 72022346388181584 |
| 9 | 4721717643249254751360 |
| 10 | 1237892809110149882059440768 |
| 11 | 1298060596773261804821355107253504 |
| 12 | 5444502293680983802677246555274553481984 |
| 13 | 91343781554246596956424128384394531707099632640 |
| 14 | 6129980884648631844901425521287946137183899295465755648 |
| 15 | 1645504465371454407878687557239154898196072267336301175996872704 |

Table 2: Number $\mathrm{C}(\mathrm{n})$ of simple labeled connected digraphs.

Simplifying Liskovets's method in 1971 Wright [100] proposed the following formulas. Let $\mathrm{n} \geq 1$,

$$
\begin{equation*}
\eta(n)=D(n)-\sum_{i=1}^{n-1} 2^{(n-1)(n-i)} \eta(i) \tag{6}
\end{equation*}
$$

and

$$
\begin{equation*}
S(n)=\eta_{n}+\sum_{i=1}^{n}\binom{n-1}{i-1} S(i) \eta_{n-i} . \tag{7}
\end{equation*}
$$

According to (6) $\eta_{1}=1, \eta_{2}=0, \eta_{3}=16$, and $\eta_{4}=1536$. Using these $\eta$ values from (7) we get $S(1)=1, S(2)=1, S(3)=18$, and $S(4)=1606$.

The values of $S(n)$ are in Table 3 for $n=1,2, \ldots, 15$. In [75] also the values for $n=16,17,18$ can be found.

In 1969 Liskovets [52] proved the following theorem.
Theorem 1 (Liskovets, 1969 [52]) If $n \geq 1$, then

$$
\begin{equation*}
D(n)-2(n+4) n^{(n+1)(n+1)} \leq S(n) \leq D(n) \tag{8}
\end{equation*}
$$

and

$$
\begin{equation*}
S(n)=D(n)\left(1-n 2^{2-n}+n(2 n-1) 2^{2-2 n}\right)+O\left(n^{3} n^{n(n-4)}\right) \tag{9}
\end{equation*}
$$

Proof. See (Liskovets, 1969 [52]).

| n | $\mathrm{S}(\mathrm{n})$ |
| ---: | ---: |
| 1 | 1 |
| 2 | 1 |
| 3 | 18 |
| 4 | 1606 |
| 5 | 565080 |
| 6 | 734774776 |
| 7 | 3523091615568 |
| 8 | 63519209389664176 |
| 9 | 11904337053178974685102609280 |
| 10 | 1270463864957828799318424676767488 |
| 11 | 95381067966826255132459611681511359329536 |
| 12 | 6109064462821545704046426032465737763224760635732888576 |
| 13 |  |
| 14 |  |

Table 3: Number $S(n)$ of simple labeled strongly connected digraph.

### 2.2 Generation of all strongly connected graphs

Let $m$ and $n$ be positive integers, $V=\left\{V_{1} \ldots, V_{n}\right\}$ be a finite set and $A=$ $\left\{a_{1}, \ldots, a_{m}\right\}$ be a finite family of ordered pairs $\left(V_{i}, V_{j}\right) \in V \times V$ of the elements of $V$. Let $D=(V, A)$ be an arbitrary directed graph [78, Volume A, page 28] and $D^{\top}=\left(V, A^{\top}\right)$ be the transpose [15, page 530, Exercise 22.1-3] of $D$ defined by

$$
\begin{equation*}
A^{\top}=\left\{\left(V_{i}, V_{j}\right) \in V \times V \mid\left(V_{j}, V_{i}\right) \in V \times V\right\} \tag{10}
\end{equation*}
$$

A directed spanning tree T of a directed graph $\mathrm{D}=(\mathrm{V}, \mathrm{A})$ is a rooted tree that consists entirely of $\operatorname{arcs}$ in $A$, all arcs directed from parents to children in the tree, and that contains every vertex of D. A directed spanning tree of D with root vertex $\mathrm{V}_{\mathrm{i}} \in \mathrm{V}$ is a breadth-first spanning tree provided that each vertex of $D$ at distance $d$ from $V_{i}$ appears at depth $d$ in the tree (that is, at distance $d$ from $V_{i}$ in the tree) [59].

We enumerated the strongly connected networks. The base of the enumeration is the fact that the strong components of a directed graph D and its transpose $\mathrm{D}^{\top}$ contain the same strongly connected components [15]. Therefore we choose arbitrary vertex as a root and build a BST (breath-first spanning tree) of the given D and of its transpose $\mathrm{D}^{\top}$. D is strongly connected if and only if both deep search trees contain all vertices of $D$.

Lemma 1 (Cormen et al., 1969 [15]) A directed graph D is strongly connected
if and only if its arbitrary vertex (e.g. $\mathrm{V}_{1}$ ) is the root of a breadth-first spanning tree of D and also the root of the breadth-first search tree $\mathrm{D}^{\top}$.

Proof. Let $V_{a}$ and $V_{b}$ be arbitrary vertices of a strongly connected graph D. Then D contains a directed path $\mathrm{V}_{\mathrm{a}}=\mathrm{V}_{\mathrm{i}_{1}}, \ldots, \mathrm{~V}_{\mathrm{i}_{\mathrm{p}}}=\mathrm{V}_{\mathrm{b}}$ and also a directed path $\mathrm{V}_{\mathrm{b}}=\mathrm{V}_{\mathrm{j}_{1}}, \ldots, \mathrm{~V}_{\mathrm{j}_{\mathrm{q}}}=\mathrm{V}_{\mathrm{a}}$. Therefore $\mathrm{D}^{\top}$ contains the directed paths $\mathrm{V}_{\mathrm{a}}=\mathrm{V}_{\mathrm{j}_{\mathrm{q}}}, \ldots, \mathrm{V}_{\mathrm{j}_{1}}=\mathrm{V}_{\mathrm{b}}$ and $\mathrm{V}_{\mathrm{b}}=\mathrm{V}_{\mathrm{i}_{\mathrm{p}}}, \ldots, \mathrm{V}_{\mathrm{i}_{1}}=\mathrm{V}_{\mathrm{a}}$, therefore the given condition is necessary.

Again let $V_{a}$ and $V_{b}$ arbitrary vertices of $D$. If $D$ contains a directed path $\left(\mathrm{V}_{1}=\mathrm{V}_{\mathrm{i}_{1}}, \ldots, \mathrm{~V}_{\mathrm{i}_{r}}=\mathrm{V}_{\mathrm{a}}\right)$ and also a directed path $\left(\mathrm{V}_{1}=\mathrm{V}_{\mathrm{j}_{1}}, \ldots, \mathrm{~V}_{\mathrm{j}_{\mathrm{q}}}=\mathrm{V}_{\mathrm{b}}\right)$, further $D^{\top}$ contains directed paths $\left(\mathrm{V}_{1}=\mathrm{V}_{\mathrm{k}_{1}}, \ldots, \mathrm{~V}_{\mathrm{k}_{\mathrm{r}}}=\mathrm{V}_{\mathrm{a}}\right)$ and $\left(\mathrm{V}_{1}=\right.$ $\left.\mathrm{V}_{\mathrm{l}_{1}}, \ldots, \mathrm{~V}_{\mathrm{l}_{\mathrm{s}}}=\mathrm{V}_{\mathrm{b}}\right)$, then D contains directed paths $\left(\mathrm{V}_{\mathrm{a}}=\mathrm{V}_{\mathrm{k}_{\mathrm{r}}}, \ldots, \mathrm{V}_{\mathrm{k}_{1}}=\right.$ $\mathrm{V}_{1}=\mathrm{V}_{\mathrm{i}_{1}}, \ldots, \mathrm{~V}_{\mathrm{i}_{\mathrm{r}}}=\mathrm{V}_{\mathrm{b}}$ ) and $\left(\mathrm{V}_{\mathrm{b}}=\mathrm{V}_{\mathrm{l}_{\mathrm{s}}}, \ldots, \mathrm{V}_{\mathrm{l}_{1}}=\mathrm{V}_{1}=\mathrm{V}_{\mathrm{i}_{1}}, \ldots, \mathrm{~V}_{\mathrm{i}_{\mathrm{r}}}=\mathrm{V}_{\mathrm{a}}\right)$, therefore the given condition is sufficient.

Algorithm Strong is based on Lemma 1. It decides if a given directed graph D is strongly connected.

Input parameters are: $\mathfrak{n}>1$ : the number of processes; $B=\left(b_{1}, \ldots, b_{n^{2}}\right)$ : the adjacency matrix of the current graph as a vector.

Output parameter is L: if D is strongly connected then $\mathrm{L}=1$, otherwise $\mathrm{L}=0$.

Working parameters are $i$ (current number of the vertices); $j, k$ : cycle variables; $m$ : the current number of vertices in the tree; $Q=\left(Q_{1}, \ldots, Q_{n}\right)$ : a queue for the waiting vertices; $h(Q)=h$ : the head index of the queue; $t(Q)=t$ : the tail index of the queue; $p=\left(p_{1}, \ldots, p_{n}\right)$ : the presence vector of the vertices $\left(p_{i}=1\right.$, if $V_{i}$ is in the tree, and $p_{i}=0$ otherwise $)$.
$\operatorname{Strong}(n, B)$
$01 \quad \mathrm{p}_{1}=1$
$\mathrm{m}=1$
$h=1$
$\mathrm{Q}_{1}=1$
$\mathrm{t}=2$
L=1
for $\mathrm{j}=2$ to n

$$
p_{j}=0
$$

while $t>h$
$/ / V_{j}$ is not in the tree.
// line 09-30: Test of D.

$$
u=Q_{h}
$$

$$
\text { for } j=1 \text { to } n-1
$$

for $k=1$ to $j-1 \quad / /$ line $12-20$ : Before the main diagonal.

13

15
16
17
if $b_{(u-1) n+k}==1$ and $p_{(u-1) n+k}==0$
// line 13: $\mathrm{V}_{\mathrm{j}}$ not in tree
$p_{k}=1 / /$ line 14-15: A new vertex of the tree is found.
$\mathrm{m}=\mathrm{m}+1$
if $m==n$
return L
$Q_{t}=j$
$\mathrm{t}=\mathrm{t}+1$
$h=h+1$
for $k=j+1$ to $n-1 / /$ line 21-30: After the main diagonal.
if $b_{(u-1) n+k}==1$ and $p_{(u-1) n+k}==0$
// line 22: $V_{j}$ not in tree.
$p_{k}=1$
$\mathrm{m}=\mathrm{m}+1$
return L
$Q_{t}=j$
$\mathrm{t}=\mathrm{t}+1$
$h=h+1$
// line 30-31: The graph is not strongly connected.

We remark that Strong tests only the existence of a breadth-first spanning tree of $D$. The test of the existence of a breadth-first spanning tree of $\mathrm{D}^{\top}$ requires similar instructions (the only difference that in lines 13 and 22 $\mathrm{b}_{(\mathfrak{u}-1) \mathfrak{n}+\mathrm{k}}=1$ must be replaced by $\mathrm{b}_{(\mathfrak{u}-1) \mathfrak{n}+\mathrm{k}}==0$.

The next assertion characterizes the resource requirements of Strong.
Theorem 2 If $\mathrm{b} \geq 2$, then Strong requires $\Theta\left(\mathrm{n}^{2}\right)$ memory locations in all cases and $\mathrm{O}\left(2^{\mathrm{b}(\mathrm{b}-1)} \mathrm{n}^{2}\right)$ time units in worst case.

Proof. The memory requirement is determined by the size of the input neighborhood matrix $B$, therefore the maximal memory requirement is $\Theta\left(n^{2}\right)$ memory locations. The time requirement of Strong is determined by the facts that the algorithm investigates at most $2^{\mathfrak{n}(n-1)}$ graphs and constructs an $\mathfrak{n} \times \mathfrak{n}$ sized matrix for all investigated graphs.

Algorithm All-Strong enumerates the strongly connected networks for $a, a+1, \ldots, b$ vertices. It is also based on Lemma 1 .
The input parameters of All-Strong are $a \geq 2$ and $b \geq a$ : lower and upper bound for the current size of the investigated network.

Output parameter is $S=(S(a), \ldots, S(b))$, where $S(a)$ is the number of the strongly connected networks consisting of a processes, ..., $S(b)$ is the number of the strongly connected networks consisting of $b$ processes.

Working parameters are $i(c u r r e n t ~ n u m b e r ~ o f ~ t h e ~ v e r t i c e s) ~ a n d ~ j ~(b o t h ~$ are cycle variables $) ; B=\left(b_{1}, \ldots, b_{n}\right)$ : the adjacency matrix of the current network as a vector; $b_{0}$ : help variable to stop the increasing of the adjacency vector; $\mathrm{Q}=\left(\mathrm{Q}_{1}, \ldots, \mathrm{Q}_{\mathrm{n}}\right)$ : a queue for the waiting vertices; $\mathrm{h}(\mathrm{Q})$ : the head index of the queue; $t(Q)=t$ : the tail index of the queue; $L$ : logical variable (if the current graph is strong, then $L=1$, otherwise $L=0$.
$\operatorname{All-Strong}(a, b)$
01 for $i=a$ to $b$
// line 01-04: Generation of the first graph.
$02 \quad S(i)=0 \quad / /$ line 02: Initialization of the enumeration.
03 for $\mathfrak{j}=0$ to $\mathfrak{i}(i-1)$
$04 \quad b_{j}=0$
$05 \operatorname{StrONG}(i, B)$
$06 \quad$ if $L==0$
07 go to 14
08 for $j=1$ to $i(i-1)$
$t_{j}=1-b_{j}$
Strong(i, T)
if $L==0 \quad / /$ line $11-12: D^{\top}$ is not strong.
go to 14
$S(i)=S(i)+1 \quad / /$ line 14: $D$ is strong
for $\mathfrak{j}=\mathfrak{i}(i-1)$ downto $1 \quad / /$ line $14-18$ : Generation of the next graph.
15
16
17
18
if $b_{j}==0$ $b_{j}=1$ for $k=j+1$ to $\mathfrak{i}(i--1)$
$b_{k}=0$
go to 05
// line 19: Continue with the next graph. print $i, S(i) \quad / /$ line 20: Print result for the current size.

The next assertion characterizes the resource requirements of All-Strong.
Theorem 3 If $\mathrm{b} \geq 2$, then All-Strong requires $\Theta(\mathrm{b}(\mathrm{b}-1)$ ) memory locations in all cases and $\mathrm{O}\left(2^{\mathrm{b}(\mathrm{b}-1)} \mathrm{n}^{2}\right)$ time units in worst case.

Proof. The memory requirement is determined by the size of the neighborhood matrices $\mathcal{B}$ and $\mathcal{T}$ defined in lines 03-04 and 08-09. The maximal size of these
matrices appears in the case when the graphs contain $b$ vertices, therefore the maximal memory requirement is $\Theta(b(b-1))$ memory locations. The time requirement of All-Strong is determined by the facts that the algorithm investigates $2^{b(b-1)}$ graphs and constructs an $n \times n$ sized matrix what according to Theorem 1 requires $\mathrm{O}\left(\mathrm{n}^{2}\right)$ time for one matrix. Multiplying these expression we get the bound $\mathrm{O}\left(2^{b(b-1)} n^{2}\right)$.

Another possible approach to generate all labeled strongly connected digraphs is to use the minimal digraphs investigated by García-López and Marijun [27].

## 3 Leader election

In the following sections the problem of leader election is considered. The mathematical models described in [59] are used: networks are modeled by directed (or sometimes undirected) graphs, processes by vertices. We suppose that the processors communicate and compute in synchronous rounds. The leader election problem is to elect a unique leader. Usually it is supposed that the processes are identical except for unique identifiers (UIDs). The size of the network is usually unknown.

In Subsection 3.1 some simple networks, then in Subsection 3.2 ring networks, in Subsection 3.3 further unidirectional networks, and finally in Subsection 3.4 further special and general networks are considered.

### 3.1 Leader election in simple networks

In this subsection the problem of leader election in simple networks as complete, chain, mesh and star networks is considered.

Peterson [71] in 1985, Afek and Gafni [1, 2] in 1981 and in 1985, Singh [81] in 1992 derived time and complexity bounds for mesh and complete networks.

In 1984 Korach et al. [47] proved optimal lower bounds for the number of messages in complete networks.

In 1985 Loui et al. [58] investigated the influence of the direction of the connections on the leader election algorithms.

There are known algorithms for chain [19] and star [80] networks too.

### 3.2 Leader election in ring networks

In this subsection comparison-based algorithms of different ring networks (in details unidirectional and bidirectional ones) are described and analyzed.

### 3.2.1 LCR algorithm in unidirectional ring

Figure 1 shows an unidirectional ring consisting of the processes $\mathrm{P}_{1}, \ldots, \mathrm{P}_{\mathrm{n}}$.


Figure 1: A ring of processes $\mathrm{P}_{1}, \ldots, \mathrm{P}_{\mathrm{n}}$.
The first known leader election algorithm was proposed by Le Lann [51] in 1977 for unidirectional rings. It is a very simple algorithm. In the first step each process sends its UID to its clockwise neighbor. In the further steps each process compares the received UID with its own UID, and if they are equal, then the process declares itself the leader, otherwise sends the larger UID to the clockwise neighbor. The algorithm terminates when the process having the largest UID gets back its own UID.

This algorithm requires $n$ steps and $n^{2}$ messages.
Chang and Roberts in 1979 [13] proposed an improved version of the previous algorithm: after the comparison of the received and own UID the processes send a message only if the received UID is the larger one. We give a formal description [59] of this algorithm called usually LCR (after Le Lann, Chang and Roberts) algorithm. It is supposed that the UID's are the natural numbers $1,2, \ldots, \mathrm{n}$.

Input parameter is n : the number of processes and $\mathrm{p}=\mathrm{p}_{1}, \ldots, \mathrm{p}_{\mathrm{n}}$ : a permutation of the UID's.

Output parameter is $M_{n}=M$ : the number of messages.
The message alphabet is $\{1,2, \ldots, n\}$. For each $i(1 \leq i \leq n)$ the state state $_{i}$ consists of three components:

- u, a UID, initially the UID of $\mathrm{P}_{\mathrm{i}}$;
- $\operatorname{send}_{i}$, a UID or null, initially the UID of $\mathrm{P}_{\mathrm{i}}$;
- statusi, having possible values \{unknown,leader $\}$.

The state of $\mathrm{P}_{\mathfrak{i}}$ consists of the single state defined by the given initial values. The message generation function $\mathrm{msgs}_{i}$ is defined by

- send the current value of send to $\mathrm{P}_{\mathrm{i}}$.

We remark that indices are interpreted everywhere mod $n$.
The transition function transi ${ }_{i}$ is defined by the following pseudocode used in [59]:

```
send := null
```

if the incoming message is v , then
case
$\mathrm{v}>\mathrm{u}$ : send $:=\mathrm{v}$
$\mathrm{v}=\mathrm{u}:$ status $_{\mathrm{i}}:=$ leader
$\mathrm{v}<\mathrm{u}$ : do nothing
endcase
Since LCR is a basic algorithm of leader election and since we execute the simulation of LCR on a sequential processor, the algorithm is described also using the pseudocode of $[15,40]$.

Input parameters are $\mathfrak{n}>1$ : the number of processes; $p=p_{1}, \ldots, p_{n}$ : a permutation of the UID's.

Output parameters are L: the index of the elected leader; $M$ : the number of messages.

Working parameters are $m=\left(m_{1}, \ldots, m_{n}\right)$, where $m_{i}$ is the current message of $\mathrm{P}_{\mathrm{i}}$; $\mathfrak{i}$ cycle variable.

## $\operatorname{LCR}(n, p)$

$01 \mathrm{P}_{\mathrm{i}}$ in parallel for $\mathfrak{i}=1$ to $n \quad / /$ line 01-05: Initialization.
$02 \quad \operatorname{read} p_{i}$
$03 \quad \mathrm{~m}_{\mathrm{i}}=\mathrm{i}$
$04 s_{i}=0$
$05 M=n$
06 while all states $s_{i}==0 \quad / /$ line 06-13: Election.
$07 \quad P_{i}$ in parallel for $\mathfrak{i}=1$ to $n$
08 if $m_{i-1}>p_{i}$
$09 \quad m_{i}=m_{i-1}$
$10 \quad M=M+1$
11 if $\mathrm{m}_{\mathrm{i}-1}=\mathrm{p}_{\mathrm{i}}$

| 12 | $s_{\mathfrak{i}}=\mathfrak{m}_{\mathfrak{i}-1}$ |
| :--- | :--- |
| 13 | $L=\mathfrak{i}$ |

14 return $L, M$ // line 14: Return of the result.
Let $X_{n}$ be a random variable characterizing the number of messages of LCR and let $M_{n}$ be the expected value of $X_{n}$ at the uniform distribution of the permutations of the UID's.

Chang and Roberts in [13] not only improved the algorithm of Le Lann, but also determined $M_{n}$.

Theorem 4 (Chang, Roberts, 1979 [13]) If the permutations of the UID's have uniform distribution, then

$$
\begin{equation*}
M_{n}=n+\sum_{i=1}^{n-1} \sum_{k=1}^{n-1} k P(n ; i, k)=n+\sum_{k=1}^{n-1} \frac{n}{k+1}=O(n \log n) \tag{11}
\end{equation*}
$$

and

$$
\begin{equation*}
M_{n}=\mathfrak{n H} H_{n}=O(n \log n) \tag{12}
\end{equation*}
$$

where $\mathrm{H}_{\mathrm{n}}$ is the n th harmonic number and $P(\mathrm{n} ; \mathrm{i}, \mathrm{k})$ is the probability that the message $\mathfrak{i}$ is passed k times.

Proof. See [13].
$P(i, k, n)$ is the probability that the $k-1$ clockwise neighbors of $i$ are less than $i$ and the kth clockwise neighbor of $i$ is larger than $i$. There are $i-1$ processes less than $\mathfrak{i}$ and $n-i$ processes larger than $i$.

Since the place of the UID $\mathfrak{i}$ can be fixed, the remaining identifiers can be permuted in $(n-1)$ ! manner. The small UID's can be choosen in $(i-1) \cdots(i-$ $k+1$ ) manner, the kth large UID in $n-i$ manner, and the remaining UID's $(n-k) \cdots 1$ manner. So we get

$$
\begin{equation*}
P(n ; i, k)=\frac{[(i-1) \cdots(i-k+1)](n-i)[(n-k) \cdots 1]}{(n-1)(n-2) \cdots 1} . \tag{13}
\end{equation*}
$$

Using the well-known bounds

$$
\begin{equation*}
\frac{1}{2}\lfloor\log n\rfloor<H_{n}<\lceil\log n\rceil \tag{14}
\end{equation*}
$$

it is easy to get the stronger assertion

$$
\begin{equation*}
M_{n}=\Theta(n \log n) \tag{15}
\end{equation*}
$$

Using Leonhard Euler's following lemma we prove Lemma 3 in which (18) and (19) are stronger than (12) in Theorem 4.

Lemma 2 (Euler [22]) If $\mathrm{n} \geq 1$ then

$$
\begin{equation*}
\mathrm{H}_{\mathrm{n}}=\sum_{\mathrm{i}=1}^{\mathrm{n}} \frac{1}{\mathfrak{i}}=\ln n+\gamma+\beta_{\mathrm{n}}, \tag{16}
\end{equation*}
$$

where $\gamma$ is the Euler-Mascheroni constant ( $\gamma \sim 0.577215$ 665) [22, 63, 96] and

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \beta_{n}=0 \tag{17}
\end{equation*}
$$

Proof. See Fichtengolz [24, Volume II, page 270].

Lemma 3 If $\mathfrak{n} \geq 1$, then

$$
\begin{equation*}
M_{n}=n \ln n+n \gamma+n \beta_{n} \tag{18}
\end{equation*}
$$

and

$$
\begin{equation*}
M_{n}=\Theta(n \log n) \tag{19}
\end{equation*}
$$

Proof. Substitution of (16) into (12) results (18) which implies (19).
Table 4 illustrates the accuracy of the approximation of (18).
Chen [14] in 2006 published a detailed probabilistic cost analysis of LCR algorithm. Using generating functions he proved

$$
\begin{equation*}
M_{n}=\frac{2}{n-1} \sum_{i=1}^{n-1} M_{i}+\frac{n}{2} \quad \text { for } \quad n \geq 2 \tag{20}
\end{equation*}
$$

and remarked that $M_{1}=1$.
Using (20) Chen reproved (11) and gave a more exact characterization

$$
\begin{equation*}
M_{n}=n \log n+\gamma n+O(1) \tag{21}
\end{equation*}
$$

of the mean of $X_{n}$, further determined the variance of $X_{n}$ as

$$
\begin{equation*}
V\left(X_{n}\right)=\left(2-\frac{\pi^{2}}{6} n^{2}\right)+O(n \log n) \tag{22}
\end{equation*}
$$

Using Euler-Maclaurin summation [21, 60, 64, 95] D. E. Knuth [46] derived the following improved version of Lemma 2.

| n | $\mathrm{E}\left(M_{\mathrm{LCR}}(\mathrm{n})\right)$ | $\mathrm{n} \ln \mathrm{n}$ | $\mathrm{n} \gamma$ | $\mathrm{n} \beta_{\mathrm{n}}$ |
| ---: | ---: | ---: | ---: | ---: |
| 1 | 1.00000000000 | 0.000000000000 | 0.5772156649015 | 0.42278433350985 |
| 2 | 3.00000000000 | 1.386294361120 | 1.154431329803 | 0.4592743090770 |
| 3 | 5.50000000000 | 3.295836866004 | 1.731646994705 | 0.4725161392911 |
| 4 | 8.33333333333 | 5.545177444480 | 2.308862659606 | 0.4792932292476 |
| 5 | 11.41666666667 | 8.047189562171 | 2.886078324508 | 0.4833987799885 |
| 6 | 14.70000000000 | 10.75055681537 | 3.463293989409 | 0.4861491952225 |
| 7 | 18.15000000000 | 13.62137104339 | 4.040509654311 | 0.4881193023021 |
| 8 | 21.74285714286 | 16.63553233344 | 4.617725319212 | 0.4895994902062 |
| 9 | 25.46071428571 | 19.77502119603 | 5.194940984114 | 0.4907521055745 |
| 10 | 29.28968253968 | 23.02585092994 | 5.772156649015 | 0.4916749607267 |
| 11 | 33.21865079365 | 26.37684800078 | 6.349372313917 | 0.4924304789518 |
| 12 | 37.23852813853 | 29.81887979746 | 6.926587978818 | 0.4930603622537 |
| 13 | 41.34173881674 | 33.34434164700 | 7.503803643720 | 0.4935935260189 |
| 14 | 45.52187257187 | 36.94680261461 | 8.081019308621 | 0.4940506486375 |
| 15 | 49.77343489843 | 40.62075301653 | 8.658234973523 | 0.4944469083788 |
| 16 | 54.09166389166 | 44.36141955584 | 9.235450638425 | 0.4947936974029 |
| 17 | 58.47239288489 | 48.16462684896 | 9.812666303326 | 0.4950997326112 |
| 18 | 62.91194540753 | 52.02669164213 | 10.38988196823 | 0.4953717971751 |
| 19 | 67.40705348573 | 55.94434060416 | 10.96709763313 | 0.4956152484385 |
| 20 | 71.95479314287 | 59.91464547108 | 11.54431329803 | 0.4958343737632 |

Table 4: Concrete values of the expressions in (18).

Lemma 4 (Knuth [46]) If $\mathrm{n} \geq 1$ then

$$
\begin{equation*}
H_{n}=\sum_{i=1}^{n} \frac{1}{\mathfrak{i}}=\ln n+\gamma+\frac{1}{2 n}+\frac{1}{12 n^{2}}+\frac{1}{120 n^{4}}-\frac{\Theta_{2, n}}{252 n^{6}} \tag{23}
\end{equation*}
$$

where $0<\Theta_{2, n}<1$.
Proof. See [46, Page 474].
It is remarkable that in the Online Encyclopedia of Integer Sequences $[83,88]$ one can find further members of the series in (23). Using the ideas of the proof of Lemma 4 we get the following characterization of $M_{n}$.

Theorem 5 If $\mathfrak{n} \geq 1$ then

$$
\begin{equation*}
M_{n}=n \ln n+\gamma n+\frac{1}{2}+\frac{1}{12 n^{2}}+\frac{1}{120 n^{4}}+\Theta\left(\frac{1}{n^{5}}\right) \tag{24}
\end{equation*}
$$

Proof. Using different methods in 1979 Chang and Roberts, in 2006 Chen proved (12). Substitution of the right side of (23) into (11) results

$$
\begin{equation*}
M_{n}=E(n)=n \ln n+\gamma n+\frac{1}{2}-\frac{1}{12 n}+\frac{1}{120 n^{3}}-\frac{\Theta_{2, n}}{252 n^{5}}, \tag{25}
\end{equation*}
$$

implying (24).
Table 5 illustrates the accuracy of the approximation of (25).

| $n$ | $\mathrm{E}(\mathrm{n})$ | $\mathrm{n} \ln n$ | $\mathrm{n} \gamma+\frac{1}{2}$ | $-\frac{1}{12 n}$ | $\frac{1}{12 n^{3}}$ | $-\frac{\Theta_{2, n}{ }^{25 n^{5}}}{}$ |
| ---: | ---: | ---: | ---: | ---: | ---: | ---: |
| 1 | 1.00000 | 0.00000 | 1.07722 | -0.0833333 | 0.0083333 | -0.0022157 |
| 2 | 3.00000 | 1.38629 | 1.65443 | -0.0416667 | 0.0010417 | -0.0001007 |
| 3 | 5.50000 | 3.29584 | 2.23165 | -0.0277778 | 0.0003086 | -0.0000147 |
| 4 | 8.33333 | 5.54518 | 2.80886 | -0.0208333 | 0.0001302 | -0.0000036 |
| 5 | 11.41667 | 8.04719 | 3.38608 | -0.0166667 | 0.0000667 | -0.0000012 |
| 6 | 14.70000 | 10.75056 | 3.96329 | -0.0138889 | 0.0000386 | -0.0000005 |
| 7 | 18.15000 | 13.62137 | 4.54051 | -0.0119048 | 0.0000243 | -0.0000002 |
| 8 | 21.74286 | 16.63553 | 5.11773 | -0.0104167 | 0.0000163 | -0.0000001 |
| 9 | 25.46071 | 19.77502 | 5.69494 | -0.0092593 | 0.0000114 | -0.0000001 |
| 10 | 29.28968 | 23.02585 | 6.27216 | -0.0083333 | 0.0000083 | -0.0000000 |
| 11 | 33.21865 | 26.37685 | 6.84937 | -0.0075758 | 0.000063 | -0.0000000 |
| 12 | 37.23853 | 29.81888 | 7.42659 | -0.0069444 | 0.0000048 | -0.0000000 |
| 13 | 41.34174 | 33.34434 | 8.00380 | -0.0064103 | 0.0000038 | -0.0000000 |
| 14 | 45.52187 | 36.94680 | 8.58102 | -0.0059524 | 0.0000030 | -0.0000000 |
| 15 | 49.77343 | 40.62075 | 9.15824 | -0.0055556 | 0.0000025 | -0.0000000 |
| 16 | 54.09166 | 44.36142 | 9.73545 | -0.0052083 | 0.0000020 | -0.0000000 |
| 17 | 58.47239 | 48.16463 | 10.31267 | -0.0049020 | 0.0000017 | -0.0000000 |
| 18 | 62.91195 | 52.02669 | 10.88988 | -0.0046296 | 0.0000014 | -0.0000000 |
| 19 | 67.40705 | 55.94434 | 11.46710 | -0.0043860 | 0.0000012 | -0.0000000 |
| 20 | 71.95479 | 59.91465 | 12.04431 | -0.0041667 | 0.0000010 | -0.0000000 |

Table 5: Concrete values of the expressions in (25).
A third possibility for the proof of (12) is the application of Pascal's next formula [68] allowing the recursive computation of the sum of the kth powers of the first $n$ positive integers.

Theorem 6 (Kovcs [49], Pascal [68], Pólya [72], Wolfram [98]) If $\mathfrak{n} \geq 1$ and
$p \geq 1$, then

$$
\begin{equation*}
S(n, p)=\sum_{i=1}^{n} \mathfrak{i}^{p}=\frac{1}{p+1}\left((n+1)^{p+1}-1-\sum_{k=1}^{p-1}\binom{p+1}{k} S(n, k)\right) . \tag{26}
\end{equation*}
$$

The following Faulhaber formula [23] also allows the computation of $S(n, p)$.
Theorem 7 (Faulhaber [23], Weisstein [97]) If $n \geq 1$ and $p \geq 1$, then

$$
\begin{equation*}
S(n, p)=\frac{1}{p+1} \sum_{i=1}^{p+1}(-1)^{\delta_{i, p}}\binom{p+1}{i} B_{p+1-i} n^{i} \tag{27}
\end{equation*}
$$

where $\delta_{i, \mathfrak{p}}$ is the Kronecker-delta [87] and $\mathrm{B}_{\mathrm{i}}$ is the Bernoulli number [84, 85].
The following double sum gives $S(n, p)$ without recursion.
Theorem 8 (Weisstein [94]) If $\mathrm{n} \geq 1$ and $\mathrm{p} \geq 1$, then

$$
\begin{equation*}
S(n, p)=\sum_{i=1}^{p} \sum_{j=0}^{i-1}(-1)^{j}(i-j)\binom{n+p-i+1}{n-i}\binom{p+1}{j} . \tag{28}
\end{equation*}
$$

### 3.2.2 Hirschberg-Sinclair algorithm in bidirectional ring

Hirschberg and Sinclair [35] in 1980 proposed an algorithm (HS) for bidirectional rings which elects as leader also the process having the largest UID. HS requires in worst case only $\Theta(n \log n)$ messages instead of the $\Theta\left(n^{2}\right)$ requirement of LCR. Figure 2 shows a bidirectional ring.

Input parameters are $n>1$ : the number of processes; $p=p_{1}, \ldots, p_{n}$ : a permutation of the UID's $1, \ldots, n$.

Output parameters: $\mathfrak{i}$ the index of the elected leader process; $\mathrm{N}=\left(\mathrm{N}_{1}, \ldots\right.$, $N_{n}$ ), where $N_{i}$ is the number of messages, sent by process $P_{i}$; $Q$ : the total number of sent messages.

Working parameters are $\mathcal{M}$ : the message alphabet $\mathfrak{m l}=\left(\mathfrak{m l}_{1}, \ldots, \mathrm{ml}_{n}\right)$, where $\mathrm{ml}_{\mathrm{i}}$ is the current message of $\mathrm{P}_{\mathrm{i}}$ to $\mathrm{P}_{\mathrm{i}-1} ; \mathrm{mr}=\left(\mathrm{mr}_{1}, \ldots, \mathrm{mr}_{n}\right)$, where $m r_{i}$ is the current message of $P_{i}$ to $P_{i+1} ; s=\left(s_{1}, \ldots, s_{n}\right)$ : status of $P_{i} ; i$ is a cycle variable; null the empty message.

The messages are triples, consisting a UID, a flag value(in or out, and a positive integer counter (hop-count) h. The possible values of the status of the processes are unknown or leader.


Figure 2: A bidirectional ring of $\mathfrak{n}$ processes.
$\operatorname{HS}(n, p)$
$01 \mathrm{P}_{\mathrm{i}}$ in parallel for $\mathfrak{i}=1$ to $n \quad / /$ line 01-05: Initialization. $02 \quad \operatorname{read} p_{i}$
$03 \quad \mathrm{ml}_{i}=(\mathrm{i}$, out, 1$) \quad / /$ line 03 : First message of $\mathrm{P}_{\mathrm{i}}$ to $\mathrm{P}_{\mathrm{i}-1}$.
$04 \quad \mathrm{mr}_{\mathfrak{i}}=(\mathfrak{i}$, out, 1$) \quad / /$ line 04: First message of $\mathrm{P}_{\mathfrak{i}}$ to $\mathrm{P}_{\mathfrak{i}+1}$.
$05 \mathrm{~s}_{\mathrm{i}}=$ unknown // line 05: Initialization of the first state of $\mathrm{P}_{\mathrm{i}}$.
$06 \mathrm{~N}=2 \mathrm{n} \quad / /$ line 06:Iinitialization of $M$.
07 while all states are unknown // line 07-12: Computation of $M$.
$07 \quad P_{i}$ in parallel for $\mathfrak{i}=1$ to $n$
$08 \quad \mathrm{mr}_{i}=$ null
$09 \quad \mathrm{ml}_{\mathrm{i}}=$ null
10 if $\mathrm{mr}_{\mathrm{i}-1}==(\mathrm{j}$, out, h$)$
$11 \quad$ if $j>i$ and $h>1$
$12 \quad \mathrm{mr}_{\mathrm{i}}=(\mathrm{j}$, out, $\mathrm{h}-1)$
$13 \quad \mathrm{~N}_{\mathrm{i}}=\mathrm{N}_{\mathrm{i}}+1$
14
15

$$
\text { if } j>i \text { and } h==1
$$

$$
\mathrm{ml}_{\mathrm{i}}=(\mathfrak{j}, i n, 1)
$$

$$
N_{i}=N_{i}+1
$$

17
18
19

$$
\text { if } \mathfrak{j}=\mathfrak{i}
$$

$$
s_{i}=\text { leader }
$$

$\mathrm{Q}=0 \quad / /$ line 17-19: Summing numbers of messages. for $i=1$ to $n$
$Q=Q+N_{i}$
return $i, N, Q \quad / /$ line 22: Return of the results.

23

$$
\text { if } \mathrm{ml}_{\mathrm{i}+1}==(\mathfrak{j}, \text { out }, \mathrm{h})
$$

if $j>i$ and $h>1$
$\mathrm{ml}_{\mathrm{i}}=(\mathrm{j}$, out, $\mathrm{h}-1)$
$\mathrm{N}_{\mathrm{i}}=\mathrm{N}_{\mathrm{i}}+1$
jf $>i$ and $h==1$
$m r_{i}=(j, i n, 1)$
$N_{i}=N_{i}+1$
$\mathbf{j} \mathbf{f}=\mathbf{i}$
$s_{i}=$ leader
$\mathrm{Q}=0 / /$ line 17-19: Summing the numbers of the messages.
for $i=1$ to $n$
$Q=Q+N_{i}$
return $i, N, Q \quad / /$ line 17: Return of the results.
if $\mathrm{ml}_{\mathrm{i}+1}==(\mathfrak{j}, i n, 1)$ and $\mathfrak{i} \neq \mathfrak{j}$
$m r_{i}=(j, i n, 1)$
$\mathrm{N}_{\mathrm{i}}=\mathrm{N}_{\mathrm{i}}+1$
if $\mathrm{ml}_{\mathrm{i}+1}==(\mathfrak{j}, i n, 1)$ and $\mathfrak{i} \neq \mathfrak{j}$
$\mathrm{ml}_{\mathrm{i}}=(\mathrm{j}, i n, 1)$
$\mathrm{N}_{\mathrm{i}}=\mathrm{N}_{\mathrm{i}}+1$
if $\mathrm{mr}_{\mathrm{i}-1}==(\mathrm{i}, i n, 1)$ and $\mathrm{ml}_{\mathrm{i}+1}==(\mathrm{i}, i n \mathbf{1})$
phase $=$ phase +1
$\mathrm{mr}_{\mathrm{i}}=\left(\mathrm{i}\right.$, out, $2^{\text {phase }}$
$\mathrm{ml}_{\mathrm{i}}=\left(\mathrm{i}\right.$, out, $2^{\text {phase }}$


Figure 3: Paths of messages of process $\mathrm{P}_{\mathfrak{i}}$ in algorithm HS.

Hirschberg and Sinclair [35] proved the following property of their algorithm. Let $W_{n}$ denote the maximal number of messages required by HS in a bidirectional synchronous ring.

Theorem 9 (Hirschberg, Sinclair [35]) If $n \geq 1$, then

$$
\begin{equation*}
W_{n} \leq 8 n(\lceil\log n\rceil+1)=\Theta(n \log n) \tag{29}
\end{equation*}
$$

and

$$
\begin{equation*}
W_{n}=O(n \log n) . \tag{30}
\end{equation*}
$$

We proved the following, stronger assertion.
Theorem 10 If $\mathrm{n} \geq 2$, then

$$
\begin{equation*}
2 n\lfloor\log n\rfloor \leq W_{n} \leq 8 n\lceil\log n\rceil \tag{31}
\end{equation*}
$$

and

$$
\begin{equation*}
W(n)=\Theta(n \log n) . \tag{32}
\end{equation*}
$$

Proof. The proof follows the ideas of application of bit reversing rings (see [59, Example 3.6.3] and [59, Figure 3.3]. Let $n=2^{k}$, for example with $k=3$. If we choose $p_{2^{0}}=p_{1}=n=8$ and $p_{2^{0}+2^{k-1}}=p_{5}=7$, then $p_{2^{0}+2^{k-2}}=p_{3}=5$, $p_{2^{0}+2^{k-1}+2^{k-2}}=p_{7}=6$, and finally the remaining processes get the UIDS $1,2,3,4$, and use similar construction for larger $k^{\prime} s$ then we need at least $8 \cdot 2+4 \cdot 2+2 \cdot 2=28$ (in general: $3,5 n$ ) messages. If $2^{k-1} \leq n<2^{k}$ then we suppose $n=2^{k}$ processes and need at least $n$ messages instead of $2 n$. If $\mathrm{n}=2$ then we need only $2 \cdot 2$ (in general: 2 n ) messages, therefore appears in the theorem only $2 n$ as lower bound.

Burns [12] published in 1980 a bidirectional algorithm which has a bit better worst case bound for the number of necessary messages.

### 3.3 Leader election is further unidirectional networks

Dolev et al. [20] and Peterson in 1982 [70] independently published an unidirectional algorithm whose worst message number is $\mathrm{O}(\mathrm{n} \log ) \mathrm{n}$, but their algorithm allows that the processes have arbitrary long response time that is they algorithm works only in asynchronous networks.

Rotem et al. [76] in 1987, Santoro et al. [77] in 1988 proposed an unidirectional asynchronous algorithm having $\mathrm{O}(\mathrm{n} \log n)$ messages in the worst case. Their algorithm elected not only the process having the largest UID, but also the processes having the k largest UID's.

Higham and Przytycka [33, 34] used a trick of Smith [89] and proposed an asynchronous algorithm what sends no more then $1.271 n \log n+O(n)$ messages in worst case.

The mentioned algorithms suppose that the processes start in the same round (otherwise they can not terminate). Recently Arrieta et al. [4] elaborated an algorithm allowing different starting rounds of the processes. The price of this property is that the guarantee for the worst message number is only $\mathrm{O}\left(\mathrm{n}^{2}\right)$.

In 1996 Alimonti et al. [3] considered the problem of choosing the minimum and maximum of the UID's when equal UID's are allowed. If the size of the ring is unknown then the problem is unsolvable. The authors describe an algorith for the unidirectional ring network containing $n$ processes, where the processes know $n$. The worst bit complexity (that is the number of sent bits) of their algorithm is $\mathrm{O}((c+\log n) n)$ with arbitrary $c>0$ and the time bound is $\mathrm{O}\left(\mathrm{c} \cdot \mathrm{n} \cdot \chi^{1 / \mathrm{c}}\right)$, where $\mathrm{x}=\max \left(\left|\mathbf{u}_{\min }\right|,\left|\mathbf{u}_{\max }\right|\right)$.

Attiya et al. [5] in 1989, Kalamboukis et al. [41] in 1991, and Pan [67] in 1994 studied the leader election in chordal rings.

Vitányi [93] in 1984 analyzed the leader election algorithms of Archimedean rings, Kranakis and Krizane [50] in 1997 of anonymous (in which the processes are undistinguishable) hypercube, and Mans [61] also in 1997 of unlabeled tori

Attiya et al. [6] proved lower bounds for the necessary number of messages for anonymous ring networks.

Ingram et al. [37] proposed a leader election algorithm for dynamic asynchronous network. Ingram et al. [36] described algorithms for dynamic networks with clausal clocks. Augustin et al. [7] published a robust leader election algorithm for the fast-changing world.

### 3.4 Leader election in further special and general networks

Peterson [71] in 1985 described efficient algorithms for mesh networks.
In 1995 Masapati and Ural [62] proposed a linear time leader election algorithms for recursively scalable networks.

Yamashita and Kameda [101], further Kranakis and Krizanc [50] investigated algorithms in anonymous hypercube networks.

Tel in 1995 [91], Flocchini and Mans [25] in 1996 analyzed the leader election algorithms of hypercube networks.

King et al. [45] in 1989, Kim and Belford [44] in 1996 proposed algorithms for unreliable networks.

In 1997 Mans [61] described an optimal distributed algorithm for unlabeled tori.

In 2001 Gavoille [29] analyzed the leader election problem of De Bruijn networks.

In 2005 Shi and Srimani [80] described an algorithm for hierarchical star
networks.
In 2007 Srimani and Lafiti [90] proposed an algorithm for Cayley networks.
In 2008 Sepehri and Godarzi [79] described an algorithm for tree networks and using heap structure they proved that their algorithm in worst case requires only $\mathrm{O}(\mathrm{n})$ messages.

Peterson [70] in 1952 described efficient algorithms for general networks.
In 1985 Afek and Gafni [2] proved that leader election in general networks requires $\Omega(n \log n)$ messages and $\Omega(\log n)$ time.

Peleg in 1990 [69] proposed a time optimal leader election algorithm for gereral networks which can be applied also for some special networks.

The basic algorithms of general networks are FloodMax and OptFloodMax (see e.g. [59]).

Das et al. [18] proposed effective algorithms which either elect a leader or signalize that the election is impossible.

Acknowledgement. The author thanks Zoltán Kása (Sapientia Hungarian University of Transylvania) for his useful critical remarks, Valery Liskovets (Mathematical Institute of Belorussian Academy of Sciences) for his help connected with the enumeration problems, István Csörgő, Attila Kovács, Sándor Kovács, and László Szili (all from Faculty of Informatics of Eötvös Loránd University) for their help connected with the sums of powers of natural numbers, PhD student Balázs Pinczel for the figures and computer experiments and PhD students Gergő Gombos and Kristóf Szabados (from the same faculty) for their technical help. The author also thanks the unknown referee for the useful remarks.

## References

[1] Y. Afek, E. Gafni, Time and message bounds for election in synchronous and asynchronous complete networks, SIAM J. Comp., 20 (1981), 376394.
[2] Y. Afek, E. Gafni, Time and message bounds for election in synchronous and asynchronous complete networks, in: Principles of Dist. Comp., ACM, 1985, 186-195.
[3] P. Alimonti, P. Flocchini, N. Santoro, Finding the extrema of a distributed multiset, J. Parallel Dist. Comp., 37 (1996), 23-33.
[4] I. Arrieta, F. Fariña, J. R. G. de Mendl, M. Raynal, Leader election: From Higham-Przytyckas algorithm to a gracefully degrading algorithm, Publications Internes de l'IRISA, inria-00605799, version 1, July 2011, 9 pages.
[5] H. Attiya, J. van Leeuwen, N. Santoro, S. Zaks, Efficient elections in chordal ring networks, Algorithmica, 4 (1989), 437-446.
[6] H. Attiya, M. Snir, M. K. Warmuth, Computing on an anonymous ring, J. ACM, 35 (1988), 845-875.
[7] J. Augustine, T. Kulkarni, P. Nakhe, P. Robinson, Robust leader election in a fast-changing world, arXive arXiv:1310.4908v1 [cs.DC], 2013, 12 pages.
[8] H. L. Bodlaender, Some lower bound results for decentralized extremafinding in rings of processors, J. Comp. System Sci., 42 (1991), 97-118.
[9] H. L. Bodlaender, New lower bound techniques for distributed leader finding and other problems on rings of processors, Theor. Comp. Sci., 81 (1991), 237-256.
[10] B. Bollobás, Random Graphs (Cambridge Studies in Advanced Mathematics 73). Cambridge University Press, Cambridge, United Kingdom, 2001, XVIII +498 pages.
[11] E. D. Burkhard, D. Kowalski, G. Malewicz, A. A. Shwartsman, Distributed algorithms, in: (ed. A. Iványi) Algorithms of Informatics, Vol. 2, mondAt Kiadó, Budapest, 2007, 591-642. Electronic version: AnTonCom, Budapest, 2011, http://progmat.hu/tananyagok/.
[12] J. E. Burns, A formal model for message passing systems. Tech. Rep. 91, Computer Science Dep., Indiana Univ., Bloomington, IN, May 1980, 21 pages.
[13] E. Chang, R. Roberts, An improved algorithm for decentralized extremafinding in circular configurations of processes, Comm. ACM, 22 (1979), 281-283.
[14] W.-M.Chen, Cost distribution of the ChangRoberts leader election algorithm, Theoret. Comp. Sci., 369 (2006), 442-447.
[15] T. H. Cormen, Ch. E. Leiserson, R. L. Rivest, C. Stein, Introduction to Algorithms (3rd edition), The MIT Press Hill, Cambridge/New York, 2009, 1312 pages.
[16] G. Coulouris, J. Dollimore, T. Kindberg, G. Blair, Distributed Systems: Concepts and Design (5th edition), Addison-Wesley, 2011, 1008 pages.
[17] G. Critzer, A recursive formula for the number of labeled simple digraphs. OEIS (ed. N. J. A. Sloane), 2012, Sequence A003027.
[18] S. Das, P. Flocchini, A. Nayak, N. Santoro, Effective elections for anonymous mobile agents, in: Algorithms and Computation, LNCS 4288, 2006, 732-743.
[19] R. Dinitz, S. Moran, S. Rajsbaum, Bit complexity of breaking and acheaving symmetry in chains and rings, J. ACM, 55 (2008), 1-28.
[20] D. Dolev, M. Klawe, M. Rodeh, An O( $n \log n$ ) unidirectional distributed algorithm for extrema finding in a circle, J. Alg., 3 (1982) 245-260.
[21] L. Euler, Methodus generalis summandi progressiones. Commentarii Academiae Scientarum Imperialis Petropolitanae, 6 (1738), 68-97, Euler Archiv E025, and Opera Omnia, 1 (1911), 42-72.
[22] L. Euler, De progressionibus harmonicus observationes, Commentarii Academiae Scientarum Imperialis Petropolitanae, 7 1740, Euler Archiv E043, 150-161 and Opera Omnia, 1 (1911), 87-100.
[23] J. Faulhaber, Academia Algebrae, Johann Remelins Verlag, Ulm, 1631, 52 pages.
[24] G. M. Fichtengolz, The Lecture on Differential and Integral Calculations, Vol. 1, 2, 3 (Russian), Nauka, Moskow, 1969, 607, 807, and 656 pages.
[25] P. Flocchini, B. Mans, Optimal elections in hypercube, J. Parallel Dist. Comp., 33 1996, 76-83.
[26] G. N. Frederickson, N. A. Lynch, Electing a leader in a synchronous ring, J. ACM, 34 1987, 98-115.
[27] J. García-López, C. Marijuán, Minimal strong digraphs, Discrete Math. 312 (2012), 737-744. Also arXiv, arXiv:1004.4827v1 [math.CO] 27 Apr 2010.
[28] H. Garcia-Molina, Election in a distributed computing system, IEEE Trans. Comp., C-31 1982, 48-59.
[29] C. Gavoille, Routing in distributed networks: Overview and open problems, ACM SIGACT News, 32 (2001), 36-52.
[30] C. Georgiou, A. A. Shvartsman, N. A. Lynch, Cooperative Task-Oriented Computing: Algorithms and Complexity (Synthesis Lectures on Distributed Computing Theory), Morgan \& Claypool Publishers, 2011, 168 pages.
[31] M. Ghaffni, N. A. Lynch, S. Sastry, Leader election using loneliness detection, in: Distributed Computing, LNCS 6950, 2011, Springer, Heidelberg, 2011, 268-282.
[32] F. Harary, Unsolved problems in the enumeration of graphs, Magyar Tud. Akad. Mat. Kutató Int. Közl. 5 (1960), 63-95.
[33] L. Higham, T. Przytycka, A simple, efficient algorithm for finding in rings, in: (ed. A. Schiper) Distributed Algorithms (7th Int. Workshp, WDAG'93, Lausanne, 1993), LNCS 725, Springer-Verlag, Berlin, 1993, 249-263.
[34] L. Higham, T. Przytycka, A simple, efficient algorithm for maximum finding on rings, Inf. Proc. Letters, 58 (1996), 319-324.
[35] D. S. Hirschberg, J. B. Sinclair, Decentralized extrema-finding in circular configuration of processes, Comm. ACM, 23 (1980), 627-628.
[36] R. Ingram, T. Radeva, P. Shields, S. Viqar, J. E. Walter, J. L. Welch, A leader election algorithm for dynamic networks with clausal clocks, Distrib. Computing, 26 (2013), 75-97.
[37] R. Ingram, P. Shields, J. E. Walter, J. L. Welch, An asynchronous leader election algorithm for dynamic networks, IEEE International Symposium on Parallel \& Distributed Processing (IPDPS 2009, Rome, 23-29 May 2009), 1-12.
[38] A. Itai, On the computational power needed to elect a leader, in: $L N C S$, 486, Springer-Verlag, Berlin, 1991, 29-40.
[39] A. Itai, M. Rodeh, Symmetry breaking in distributed networks, Inf. Comp. 88 (1990), 60-87.
[40] A. Iványi, Parallel Algorithms (Hungarian), ELTE Eötvös Kiadó, Budapest, 2005.
[41] T. Z. Kalamboukis, S. L. Mantzaris, Towards optimal distributed election on chordal rings, Information Proc. Letters, 38 (1991), 265-270.
[42] M. Kalpathi, H. M. Mahmoud, M. D. Ward, Asymptotic properties of a leader election algorithm, J. Appl. Probab., 48 (2011), 569-575.
[43] A. D. Kshemkalyani, M. Singhal, Distributed Computing: Principles, Algorithms, and Systems, Cambridge University Press, Cambridge, 2011, 756 pages.
[44] G. Kim, G. Belford, A distributed election protocol for unreliable networks, J. Parallel Distr. Comp., 35 (1996), 35-42.
[45] C.-T. King, T. B. Gendreau, L. M. Ni, Reliable election in broadcast networks, J. Parallel Distr. Computing, 7 (1989), 521-540.
[46] D. E. Knuth, Concrete Mathematics (2nd edition), Addison-Wesley Publishing Co., 1994, 672 pages (1st edition: 1988).
[47] E. S. Korach, S. Moran, S. Zaks, Optimal lower bounds for some distributed algorithms for a complete network of processors, Theoretical Comp. Sci., 64 (1989), 125-132.
[48] A. Kovács, Sums of the kth powers for the first twenty positive integers. Manuscript, Budapest, 2012.
[49] S. Kovács, Zur Berechnung der Potenzsummen. Manuscript. Budapest, 2013, 11 pages.
[50] E. Kranakis, D. Krizanc, Distributed computing on anonymous hypercube networks, J. Alg., 23 (1997), 32-50.
[51] G. Le Lann, Distributed systems-towards a formal approach, in: (ed. B. Gilricst) Information Processing 77 (Toronto, 1977).Vol. 7. of Proc. of IFIP Congress, North Holland, Amsterdam, 1977, 155-160.
[52] V. A. Liskovets, On a recurrent method for enumeration of graphs with labeled vertices (Russian), Dokl. AN SSSR, 184 (1969), 1284-1287.
[53] V. A. Liskovets, The number of strongly connected oriented graphs (Russian), Mat. Zametki, 8 (1970), 721-732.
[54] V. A. Liskovets, A contribution to the enumeration of strongly connected digraphs (Russian), Dokl. AN BSSR, 17 (1973), 1077-1080, 1163.
[55] V. A. Liskovets, On a general enumerative scheme for labeled graphs (Russian), Dokl. AN BSSR, 21 (1977), 496-499.
[56] V. A. Liskovets, Some easily derivable integer sequences. J. Int. Sequences, 3 (2000), Article 00.2.2.
[57] V. A. Liskovets, Exact enumeration of acyclic deterministic automata. Discrete Appl. Math., 154, (2006), 537-551.
[58] M. C. Loui, T. A. Matsushita, D. B. West, Election in a complete network with a sense of direction, Inform. Proc. Letters, 22 (1985), 185-187.
[59] N. A. Lynch, Distributed Algorithms (5th edition, The Morgan Kaufmann Series in Data Management Systems), Morgan Kaufmann Publishers, 2003, XIII +873 pages (1st edition: 1996).
[60] C. MacLaurin, A Treatise of Fluxions, Vol. 1. and 2, T. W. Ruddimans and T. Ruddimans, Edinburgh, 1742, 763 pages.
[61] B. Mans, Optimal distributed algorithms in unlabeled tori and chordal rings, J. Parallel Distributed Comp., 46 (1997), 80-90.
[62] G. H. Masapati, H. Ural, Electing a leader in a synchronous recursively scalable network, in ICCI90, LNCS 468, Springer-Verlag, Berlin, 1990, 463-472.
[63] L. Mascheroni, Ad notationes ad calculum integralem Euleri, Vol. 1 and 2. Ticino, Italy, 1790 and 1792. Reprinted in Euler, L. Leonhardi Euleri Opera Omnia, Ser. 1, Vol. 12, Teubner, Leipzig, Germany, 415-542.
[64] Yu. V. Matiyasevich, Alternatives to the Euler-Maclaurin formula for calculating infinite sums, Math. Notes, 88 (2010), 524-529.
[65] T. D. Noe, Number of labeled weakly connected digraphs with n nodes for $\mathrm{n}=1, \ldots, 35$. In OEIS (ed. N. J. A. Sloane), May 11, 2007. http://oeis.org/A053763/b053763.txt.
[66] T. D. Noe, Number of labeled simple connected digraphs with $\mathfrak{n}$ nodes for $\mathrm{n}=1, \ldots, 30$. In OEIS (ed. N. J. A. Sloane), January 9, 2009. http://oeis.org/A003027/b003027.txt.
[67] Y. Pan, A near-optimal multistage distributed algorithm for finding leaders in clustered chordal rings, Information Sci., 76 (1994), 131-140.
[68] B. Pascal, Ouvres de Blaise Pascal, Vol. 3 (ed L. Brunschvicg, P. Bourroux), Nabu Press, Charleston, SC, 2010, 341-367. 1st edition: Blaise Pascal, Ouvres, 1640.
[69] D. Peleg, Time-optimal leader election in general networks, J. Parallel Distr. Comp., 8 (1990), 96-99.
[70] G. L. Peterson, $\mathrm{An} \mathrm{O}(\mathrm{n} \log \mathrm{n})$ unidirectional distributed algorithm for the circular extremal problem, ACM Trans. Lang. Systems, 4 (1982), 758-762.
[71] G. L. Peterson, Efficient algorithms for elections in meshes and complete networks, TR 140, Dept. of Computer Science, Univ. of Rochester, 1985, 5 pages.
[72] Gy. Pólya, Mathematical Discovery on Understanding, Learning and Teaching Problem Solving. John Wiley \& Sons, Inc, New York, NY, 1962, 216 pages.
[73] R. W. Robinson, Counting labeled acyclic digraphs, in: New Directions in the Theory of Graphs (F. Harary, ed.), Academic Press, New York, 1973, 239-273.
[74] R. W. Robinson, Counting unlabeled acyclic digraphs, Combinatorial Mathematics V. Lecture Notes in Math., 622 (1977), 28-43.
[75] R. W. Robinson, Table of $\mathfrak{n}, \mathfrak{a}(\mathfrak{n})$ for $\mathfrak{n}=1, \ldots, 18$, in OEIS (ed. N. J. A. Sloane), 2012. Sequence A003030.
[76] D. Rotem, E. Korach, N. Santoro, Analysis of a distributed algorithm for extrema finding in a ring, J. Parallel Distr. Comp., 4 (1987), 575-591.
[77] N. Santoro, M. Scheutzow, J. B. Sidney, On the expected complexity of distributed selection, J. Parallel Distr. Comp., 5 (1988), 194-203.
[78] A. Schrijver, Combinatorial Optimization. Vol. A, B, C (Algorithms and Combinatorics, Vol. 24, Springer-Verlag, Berlin, 2003, 1800 pages.
[79] M. Seperhi, M. Godarzi, Leader election algorithm using heap structure, in: 12th WSEAS Int.l Conf. on Computers (Heraklion, Greece, July 2325, 2008), 2008, 668-672.
[80] W. Shi, K. Srimani, Leader election in hierarchical star network, J. Parallel Distr. Comp., 65 (2005), 1435-1442.
[81] G. Singh, Leader election in complete networks, in: Proc. of the Eleventh Annual ACM Symp. on Principles of Distributed Computing, ACM Press, 1992, 179-190.
[82] N. J. A. Sloane, Number of directed graphs (or digraphs) with n nodes. OEIS (ed. N. J. A. Sloane), 2013, Sequence A000273.
[83] N. J. A. Sloane, Number of strongly connected digraphs with n labeled nodes, OEIS (ed. N. J. A. Sloane), 2013, Sequence A003030.
[84] N. J. A. Sloane, Numerator of Bernoulli number B ${ }_{n}$. OEIS (ed. N. J. A. Sloane), 2013, Sequence A027641.
[85] N. J. A. Sloane, Denominator of Bernoulli number B ${ }_{n}$., OEIS (ed. N. J. A. Sloane), 2013, Sequence A027642.
[86] N. J. A. Sloane, The number of directed graphs on $\mathfrak{n}$ vertices, OEIS (ed. N. J. A. Sloane), 2013, Sequence A003085.
[87] N. J. A. Sloane, Jacobi (or Knonecker) symbol, OEIS (ed. N. J. A. Sloane), 2013, Sequence A034947.
[88] N. J. A. Sloane, Bernoulli Numbers B $2 n / 2 n$, in OEIS (ed. N. J. A. Sloane), 2013, Sequence A006953.
[89] A. R. Smith, Cellular automata complexity trade-offs, Inf. Control., 18 (1971), 466-482.
[90] P. Srimani, S. Lafiti, Some bounded degree communication networks and optimal leader election, in: Combinatorial Optimization in Communication Networks (Combinatorial Optimization), 18 (2006), 467-501.
[91] G. Tel, Linear election in hypercubes, Parallel Proc. Letters, 5 (1995), 357-366.
[92] G. Tel, Introduction to Distributed Systems (Second edition), Cambridge University Press, Cambridge, 2000, 612 pages (1st edition appeared in 1984).
[93] P. Vitányi, Distributed elections in archimedean ring of processors, in Proc. 16th Ann. ACM Symp. on Theory of Computing, 1984, 542-547.
[94] E. W. Weisstein, Double Sum, From Mathworld-A Wolfram Web Resource, 2013, http://mathworld.wolfram.com/PowerSum.html.
[95] E. W. Weisstein, Euler-Maclaurin Integration Formulas, 2013, http://mathworld.wolfram.com/Euler-MaclaurinIntegrationFormulas.html.
[96] E. W. Weisstein, Euler-Mascheroni Constant, From Mathworld-A Wolfram Web Resource, 2013, http://mathworld.wolfram.com/Euler-MascheroniConstant.html.
[97] E. W. Weisstein, Power Sum, From Mathworld-A Wolfram Web Resource, 2013, http://mathworld.wolfram.com/PowerSum.html.
[98] S. Wolfram, Wolframalpha, 2013. http://www.wolframalpha.com.
[99] E. M. Wright, Asymptotic enumeration of connected graphs. Proc. Roy. Soc. Edinburgh Sect. A, 68 (1968/1970), 298-308.
[100] E. M. Wright, The number of strong digraphs, Bull. London Math. Soc., 3 (1971), 348-350.
[101] M. Yamashita, T. Kameda, Computing on anonymous networks: Parts I and II, IEEE Trans. Par. Dist. Syst., 7 (1996), 69-96.

# About a condition for starlikeness 

Róbert Szász<br>Sapientia Hungarian University of Transylvania<br>Department of Mathematics and Informatics<br>Târgu Mureş, Romania<br>email: rszasz@ms.sapientia.ro

Pál Kupán<br>Sapientia Hungarian University of<br>Transylvania<br>Department of Mathematics and<br>Informatics<br>Târgu Mureş, Romania<br>email: kupanp@ms.sapientia.ro

## Dedicated to the memory of Professor Antal Bege


#### Abstract

A result concerning the starlikeness of the image of the Alexander operator is improved in this paper. The techniques of differential subordinations are used.


## 1 Introduction

Let $\mathbb{U}\left(z_{0}, r\right)=\left\{z \in \mathbb{C}| | z-z_{0} \mid<r\right\}$ be the disk centred in $z_{0}$ and let $U=U(0,1)$ be the open unit disk in $\mathbb{C}$. Let $\mathcal{A}$ be the class of analytic functions f , which are defined on the unit disc $U$ and have the form: $f(z)=z+a_{2} z^{2}+a_{3} z^{3}+\cdots$.

The subclass of $\mathcal{A}$ consisting of functions for which the domain $f(U)$ is starlike with respect to 0 , is denoted by $S^{*}$. An analytic characterization of $S^{*}$ is given by

$$
S^{*}=\left\{\mathrm{f} \in \mathcal{A}: \mathfrak{R} \frac{z \mathrm{f}^{\prime}(z)}{\mathrm{f}(z)}>0, z \in \mathrm{U}\right\}
$$

Another subclass of $\mathcal{A}$ we deal with is the class of close-to-convex functions denoted by $C$. A function $\mathrm{f} \in \mathcal{A}$ belongs to the class C if and only if there is a starlike function $g \in S^{*}$, so that $\mathfrak{R} \frac{z f^{\prime}(z)}{g(z)}>0, \quad z \in U$. We note that $C$ and

## 2010 Mathematics Subject Classification: 30C45

Key words and phrases: Alexander operator, starlike functions, close-to-convex functions

S* contain univalent functions. The Alexander integral operator is defined by the equality

$$
A(f)(z)=\int_{0}^{z} \frac{f(t)}{t} d t
$$

The authors of [1] (p. 310-311) proved the following result:
Theorem 1 Let $\mathcal{A}$ be Alexander operator and let $\mathrm{g} \in \mathcal{A}$ satisfy

$$
\begin{equation*}
\mathfrak{R} \frac{z g^{\prime}(z)}{g(z)} \geq\left|\mathfrak{I} \frac{z\left(z g^{\prime}(z)\right)^{\prime}}{g(z)}\right|, z \in \mathrm{U} . \tag{1}
\end{equation*}
$$

If $\mathrm{f} \in \mathcal{A}$ and

$$
\begin{equation*}
\mathfrak{R} \frac{z f^{\prime}(z)}{g(z)}>0, z \in U \tag{2}
\end{equation*}
$$

or

$$
\begin{equation*}
\mathfrak{R} \frac{f^{\prime}(z)}{g^{\prime}(z)}>0, z \in U \tag{3}
\end{equation*}
$$

then $\mathrm{F}=\mathrm{A}(\mathrm{f}) \in \mathrm{S}^{*}$.
In [1], [3], [5] improvements of the first part ((1), (2) $\left.\Rightarrow A(f) \in S^{*}\right)$ of this result is proved, simplifying condition (1). The aim of this paper is to give an improvement for the second part of Theorem 1. In order to do this, we need the definitions and lemmas exposed in the next section.

## 2 Preliminaries

Let $f$ and $g$ be analytic functions in $U$. The function $f$ is said to be subordinate to $g$, written $f \prec g$, if there is a function $w$ analytic in $U$, with $w(0)=0$, $|w(z)|<1, z \in U$ and $f(z)=g(w(z)), z \in U$. Recall that if $g$ is univalent, then $f \prec g$ if and only if $f(0)=g(0)$ and $f(U) \subset g(U)$.

Lemma 1 [2] p. 24 (Miller-Mocanu)
Let $\mathrm{p}(z)=\mathrm{a}+\sum_{\mathrm{k}=\mathrm{n}}^{\infty} \mathrm{a}_{\mathrm{k}} z^{\mathrm{k}}$ be analytic in U with $\mathrm{p}(z) \not \equiv \mathrm{a}, \mathrm{n} \geq 1$ and let $\mathrm{q}: \mathrm{U} \rightarrow \mathbb{C}$ be an analytic and univalent function with $\mathrm{q}(0)=\mathrm{a}$. If p is not subordinate to q , then there are two points $z_{0} \in \mathrm{U},\left|z_{0}\right|=\mathrm{r}_{0}$ and $\zeta_{0} \in \partial \mathrm{U}$ and a real number $\mathrm{m} \in[\mathrm{n}, \infty)$, so that q is defined in $\zeta_{0}$, $\mathrm{p}\left(\mathrm{U}\left(0, \mathrm{r}_{0}\right)\right) \subset \mathrm{q}(\mathrm{U})$, and:
(i) $\mathrm{p}\left(z_{0}\right)=\mathrm{q}\left(\zeta_{0}\right)$,
(ii) $z_{0} p^{\prime}\left(z_{0}\right)=m \zeta_{0} q^{\prime}\left(\zeta_{0}\right)$,
and
(iii) $\operatorname{Re}\left(1+\frac{z_{0} \mathrm{p}^{\prime \prime}\left(z_{0}\right)}{\mathrm{p}^{\prime}\left(z_{0}\right)}\right) \geq m \operatorname{Re}\left(1+\frac{\zeta_{0} \mathrm{q}^{\prime \prime}\left(\zeta_{0}\right)_{0}}{\mathrm{q}^{\prime}\left(\zeta_{0}\right)}\right)$.

We note that $z_{0} \mathrm{p}^{\prime}\left(z_{0}\right)$ is the outward normal to the curve $\mathrm{p}\left(\partial \mathrm{U}\left(0, \mathrm{r}_{0}\right)\right)$ at the point $\mathfrak{p}\left(z_{0}\right)$. $\left(\partial \mathrm{U}\left(0, \mathrm{r}_{0}\right)\right.$ denotes the border of the disc $\mathrm{U}\left(0, \mathrm{r}_{0}\right)$.)

Lemma 2 [2] p. 26 (Miller-Mocanu) Let $p(z)=a+\sum_{k=n}^{\infty} a_{k} z^{k}, p(z) \not \equiv a$ and $n \geq 1$.
If $z_{0} \in \mathrm{U}$ and

$$
\operatorname{Rep}\left(z_{0}\right)=\min \left\{\operatorname{Rep}(z):|z| \leq\left|z_{0}\right|\right\},
$$

then
(i) $z_{0} p^{\prime}\left(z_{0}\right) \leq-\frac{n}{2} \frac{\left|p\left(z_{0}\right)-a\right|^{2}}{\operatorname{Re}\left(a-p\left(z_{0}\right)\right)}$
and
(ii) $\operatorname{Re}\left[z_{0}^{2} \mathrm{p}^{\prime \prime}\left(z_{0}\right)\right]+z_{0} \mathrm{p}^{\prime}\left(z_{0}\right) \leq 0$.

Lemma 3 If $\mathrm{d}=\frac{2}{\pi} \arctan \left(\frac{1}{2.273}\right)$, and $\mathrm{k}_{\mathrm{d}}(z)=\int_{0}^{z} \frac{\left(\frac{1+\mathrm{t}}{}\right)^{\mathrm{d}} 1-\mathrm{t}}{\mathrm{t}} \mathrm{dt}$, then

$$
\left|\Im\left(k_{\mathrm{d}}(z)\right)\right| \leq \frac{\pi}{6}, \quad z \in \mathrm{U} .
$$

Proof. The maximum principle for harmonic functions implies that

$$
\sup _{z \in \mathrm{U}}\left|\mathfrak{J} \mathrm{k}_{\mathrm{d}}(z)\right|=\sup _{\theta \in[-\pi, \pi]}\left|\mathfrak{I} \mathrm{k}_{\mathrm{d}}\left(e^{\mathfrak{i \theta}}\right)\right| .
$$

On the other hand we have:

$$
v_{\mathrm{d}}(\theta)=\mathfrak{J k _ { \mathrm { d } }}\left(e^{\mathrm{i} \theta}\right)=\int_{0}^{1} \frac{1}{x}\left(\frac{1+x^{2}+2 x \cos \theta}{1+x^{2}-2 x \cos \theta}\right)^{\mathrm{d}} \sin \left(\mathrm{~d} \arctan \left(\frac{2 x \sin \theta}{1-x^{2}}\right)\right) \mathrm{dx} .
$$

This implies that $v_{\mathrm{n}}$ is an even function, consequently

$$
\sup _{\theta \in[-\pi, \pi]}\left|\Im k_{\mathrm{d}}\left(e^{\mathrm{i} \theta}\right)\right|=\sup _{\theta \in[0, \pi]}\left|\Im k_{\mathrm{d}}\left(e^{\mathrm{i} \theta}\right)\right| .
$$

We will prove the following equality:

$$
\begin{array}{r}
\mathrm{k}_{\mathrm{d}}\left(e^{\mathrm{i} \mathrm{\theta}}\right)=\int_{0}^{1} \frac{\left(\frac{1+x e^{i \theta}}{1-x e^{i \theta}}\right)^{\mathrm{d}}-1}{x} \mathrm{~d} x=\int_{0}^{\infty}\left[\left(\frac{e^{\mathrm{t}}-1}{e^{\mathrm{t}}+1}\right)^{\mathrm{d}}-1\right] \mathrm{dt}+\mathfrak{i}(\pi-\theta)+ \\
\quad\left(\sin \left(\frac{\pi}{2} \mathrm{~d}\right)-i \cos \left(\frac{\pi}{2} \mathrm{~d}\right)\right) \int_{0}^{\pi-\theta} \tan ^{\mathrm{d}} \frac{x}{2} \mathrm{~d} x, \theta \in[0, \pi] . \tag{4}
\end{array}
$$

We begin with the observation that the change of variable $x=e^{-t}$ leads to

$$
k_{\mathrm{d}}\left(e^{\mathfrak{i} \theta}\right)=\int_{0}^{\infty}\left[\left(\frac{e^{\mathrm{t}}+e^{i \theta}}{e^{\mathrm{t}}-e^{\mathrm{i} \mathrm{\theta}}}\right)^{\mathrm{d}}-1\right] \mathrm{dt} .
$$

Let $\theta \in[0, \pi]$ and consider the function

$$
f(z)=\left(\frac{e^{z}+e^{i \theta}}{e^{z}-e^{i \theta}}\right)^{d}-1 .
$$

We integrate it on $\Gamma=\gamma_{1} \cup \gamma_{2} \cup \gamma_{3} \cup \gamma_{4}$, where $\gamma_{1}(t)=t, t \in[0, R], \gamma_{2}(t)=$ $\mathrm{R}-\mathrm{it}, \mathrm{t} \in[0, \pi-\theta], \gamma_{3}(\mathrm{t})=\mathrm{R}-\mathrm{t}+\mathfrak{i}(\theta-\pi), \mathrm{t} \in[0, \mathrm{R}]$ and $\gamma_{4}(\mathrm{t})=$ $\mathfrak{i}(\theta-\pi+t), t \in[0, \pi-\theta]$. The obtained equality $\int_{\Gamma} f(z) d z=0$ leads to

$$
\begin{gathered}
k_{d}\left(e^{i \theta}\right)=\lim _{R \rightarrow \infty} \int_{\gamma_{1}} f(z) d z=-\lim _{R \rightarrow \infty}\left[\int_{\gamma_{2}} f(z) d z+\int_{\gamma_{3}} f(z) d z+\right. \\
\left.\int_{\gamma_{4}} f(z) d z\right]=\int_{0}^{\infty}\left[\left(\frac{e^{x}-1}{e^{x}+1}\right)^{d}-1\right] d x+ \\
\mathfrak{i}(\pi-\theta)+\left(\sin \left(\frac{\pi}{2} d\right)-\mathfrak{i} \cos \left(\frac{\pi}{2} d\right)\right) \int_{0}^{\pi-\theta} \tan ^{d} \frac{x}{2} d x .
\end{gathered}
$$

Thus, it follows

$$
v_{\mathrm{d}}(\theta)=\mathfrak{I k}_{\mathrm{d}}\left(e^{\mathrm{i} \theta}\right)=\pi-\theta-\cos \left(\frac{\pi}{2} \mathrm{~d}\right) \int_{0}^{\pi-\theta} \tan ^{\mathrm{d}} \frac{x}{2} \mathrm{~d} x .
$$

The function $v_{\mathrm{d}}:[0, \pi] \rightarrow \mathbb{R}$ has a maximum at the point $\theta_{\mathrm{d}}=2 \arctan$ $\left(\cos ^{\frac{1}{d}}\left(\frac{\pi}{2} \mathrm{~d}\right)\right)$. A suitable numerical approach shows that

$$
\left|\mathfrak{I}\left(k_{d}(z)\right)\right| \leq v_{d}\left(\theta_{d}\right)=0.49 \cdots<\frac{\pi}{6}
$$

Lemma 4 If $\mathrm{q}_{\mathrm{d}}(z)=\exp \left(\int_{0}^{z} \frac{\left(\frac{1+\mathrm{t}}{1-\mathrm{t}}\right)^{\mathrm{d}}-1}{\mathrm{t}} \mathrm{dt}\right)=\exp \left(\mathrm{k}_{\mathrm{d}}(z)\right), \mathrm{p} \in \mathcal{A}$, and

$$
\frac{z \mathrm{p}^{\prime}(z)}{\mathrm{p}(z)} \prec \mathrm{h}(z)=\frac{z \mathrm{q}_{\mathrm{d}}^{\prime}(z)}{\mathrm{q}_{\mathrm{d}}(z)}, z \in \mathrm{U},
$$

then $p \prec q_{d}$.

Proof. We have: $\frac{1}{2 \mathrm{~d}} \mathrm{~h}(z)=\frac{1}{2 \mathrm{~d}}\left(\left(\frac{1+\mathrm{t}}{1-\mathrm{t}}\right)^{\mathrm{d}}-1\right) \in \mathrm{S}^{*}$. Lemma 3 implies:

$$
\mathfrak{R} \exp \left(\int_{0}^{z} \frac{\left(\frac{1+\mathrm{t}}{1-\mathrm{t}}\right)^{\mathrm{d}}-1}{\mathrm{t}} \mathrm{dt}\right)>0, z \in \mathrm{U} .
$$

On the other hand:

$$
\frac{z q_{d}^{\prime}(z)}{h(z)}=\exp \left(\int_{0}^{z} \frac{\left(\frac{1+t}{1-t}\right)^{d}-1}{t} d t\right)
$$

These imply $\mathrm{q}_{\mathrm{d}} \in \mathrm{C}$, which means that $\mathrm{q}_{\mathrm{d}}$ is univalent. If the subordination $p \prec q_{d}$ does not hold, then according to the Miller-Mocanu lemma it follows that there are two points, $z_{0} \in \mathrm{U}$ and $\zeta_{0} \in \partial \mathrm{U}$, and a real number $\mathrm{m} \in[1, \infty)$ such that

$$
\begin{aligned}
p\left(z_{0}\right) & =q_{d}\left(\zeta_{0}\right), \\
z_{0} p^{\prime}\left(z_{0}\right) & =m \zeta_{0} q_{d}^{\prime}\left(\zeta_{0}\right) .
\end{aligned}
$$

Since $h(U)$ is a starlike domain with respect to 0 , it follows that:

$$
\frac{z_{0} p^{\prime}\left(z_{0}\right)}{p\left(z_{0}\right)}=m \frac{\zeta_{0} q_{d}^{\prime}\left(\zeta_{0}\right)}{q_{\mathrm{d}}\left(\zeta_{0}\right)}=m h\left(\zeta_{0}\right) \notin \mathrm{h}(\mathrm{U}) .
$$

This contradicts the subordination $\frac{\operatorname{zp}^{\prime}(z)}{\mathfrak{p}(z)} \prec h(z), z \in U$. The obtained contradiction implies: $\mathrm{p} \prec \mathrm{q}_{\mathrm{d}}$.

Lemma 5 If $\mathrm{f} \in \mathcal{A}$ and

$$
\left|\arg \frac{z g^{\prime}(z)}{\mathrm{g}(z)}\right|<\arctan \left(\frac{1}{2.273}\right), z \in \mathrm{U}
$$

then

$$
\left|\arg \frac{\mathrm{g}(z)}{z}\right|<\frac{\pi}{6}, \quad z \in \mathrm{U}
$$

Proof. The condition of the lemma is equivalent to

$$
\frac{z g^{\prime}(z)}{\mathrm{g}(z)} \prec\left(\frac{1+z}{1-z}\right)^{\mathrm{d}}, z \in \mathrm{u} .
$$

Replacing in the previous lemma $p(z)=\frac{\mathrm{g}(z)}{z}$, we get

$$
\frac{\mathrm{g}(z)}{z} \prec \mathrm{q}_{\mathrm{d}}(z)=\exp \left(\int_{0}^{z} \frac{\left(\frac{1+\mathrm{t}}{1-\mathrm{t}}\right)^{\mathrm{d}}-1}{\mathrm{t}} \mathrm{dt}\right), \quad z \in \mathrm{U} .
$$

Thus

$$
\left|\arg \frac{\mathrm{g}(z)}{z}\right| \leq \max _{\theta \in[-\pi, \pi]}\left|\mathfrak{I} \int_{0}^{1} \frac{\left(\frac{1+e^{i \theta} \mathrm{t}}{1-e^{i \theta} t}\right)^{\mathrm{d}}-1}{\mathrm{t}} \mathrm{dt}\right|=v_{\mathrm{d}}\left(\theta_{\mathrm{d}}\right)<\frac{\pi}{6}, \quad z \in \mathrm{U}
$$

In [1] the following theorem is proved.
Theorem 2 If $\mathrm{f} \in \mathcal{A}$, end

$$
\begin{equation*}
\mathfrak{R} \frac{z g^{\prime}(z)}{\mathrm{g}(z)} \geq\left|\mathfrak{I} \frac{z\left(z \mathrm{~g}^{\prime}(z)\right)^{\prime}}{\mathrm{g}(z)}\right|, \quad z \in \mathrm{U} \tag{5}
\end{equation*}
$$

then the following inequality holds:

$$
\begin{equation*}
\mathfrak{R} \frac{z \mathrm{~g}^{\prime}(z)}{\mathrm{g}(z)}>2.273\left|\mathfrak{I} \frac{z \mathrm{~g}^{\prime}(z)}{\mathrm{g}(z)}\right|, \quad z \in \mathrm{U} \tag{6}
\end{equation*}
$$

## 3 The main result

Theorem 3 If $\mathrm{g} \in \mathcal{A}$ satisfies (5), then

$$
\begin{equation*}
\left|\arg \left(g^{\prime}(z)\right)\right|<\frac{5 \pi}{17}, \quad z \in \mathrm{U} \tag{7}
\end{equation*}
$$

Proof. Inequality (6) is equivalent to

$$
\begin{equation*}
\left|\arg \frac{z \mathrm{~g}^{\prime}(z)}{\mathrm{g}(z)}\right|<\arctan \frac{1}{2.273}, \quad z \in \mathrm{U} \tag{8}
\end{equation*}
$$

Thus according to Lemma 5 the inequality

$$
\left|\arg \frac{g(z)}{z}\right|<\frac{\pi}{6}, \quad z \in U
$$

follows. Summarizing we get

$$
\left|\arg g^{\prime}(z)\right| \leq\left|\arg \frac{z g^{\prime}(z)}{g(z)}\right|+\left|\arg \frac{g(z)}{z}\right|<\arctan \frac{1}{2.273}+\frac{\pi}{6}<0.92<\frac{5 \pi}{17}
$$

If we could improve the previously proved result proving that $\left|\arg \left(g^{\prime}(z)\right)\right|<$ $\frac{\pi}{5}, z \in \mathrm{U}$, then it would follow that the next theorem is an improvement of Theorem 1.

Theorem 4 If $\mathrm{f}, \mathrm{g} \in \mathcal{A}$ and

$$
\begin{equation*}
\left|\arg \left(g^{\prime}(z)\right)\right|<\frac{\pi}{5}, z \in \mathrm{U}, \tag{9}
\end{equation*}
$$

then the condition

$$
\mathfrak{R} \frac{f^{\prime}(z)}{g^{\prime}(z)}>0, z \in U
$$

implies that $\mathrm{F}=\mathrm{A}(\mathrm{f}) \in \mathrm{S}^{*}$.
Proof. The conditions of the theorem imply

$$
\begin{equation*}
\left|\arg f^{\prime}(z)\right| \leq\left|\arg \frac{f^{\prime}(z)}{g^{\prime}(z)}\right|+\left|\arg g^{\prime}(z)\right| \leq \frac{7 \pi}{10}, z \in \mathrm{U} . \tag{10}
\end{equation*}
$$

Using this result, we will prove that

$$
\begin{equation*}
\left|\arg \frac{f(z)}{z}\right| \leq \alpha_{0}=\frac{50 \pi}{108}, z \in \mathrm{U} . \tag{11}
\end{equation*}
$$

To do this we rewrite inequality (11) in the following equivalent form:

$$
\frac{\mathrm{f}(z)}{z} \prec\left(\frac{1+z}{1-z}\right)^{\frac{2}{\pi} \alpha_{0}}, z \in \mathrm{U} .
$$

If this subordination does not hold, then using Lemma 1 it follows that there are two points $z_{0} \in U, \zeta_{0}=e^{i \theta_{0}} \in \partial U$ and a real number $m_{0} \in[1, \infty)$, such that:

$$
\begin{aligned}
& \frac{f\left(z_{0}\right)}{z_{0}}=\left(\frac{1+\zeta_{0}}{1-\zeta_{0}}\right)^{\frac{2}{\pi} \alpha_{0}}=\left(i \cot \frac{\theta_{0}}{2}\right)^{\frac{2}{\pi} \alpha_{0}} \\
\left.z\left(\frac{f(z)}{z}\right)^{\prime}\right|_{z=z_{0}} & =f^{\prime}\left(z_{0}\right)-\frac{f\left(z_{0}\right)}{z_{0}}=\frac{2}{\pi} \mathfrak{m}_{0} \alpha_{0} \zeta_{0}\left(\frac{1+\zeta_{0}}{1-\zeta_{0}}\right)^{\frac{2}{\pi} \alpha_{0}-1} \frac{2}{\left(1-\zeta_{0}\right)^{2}} \\
& =\frac{2}{\pi} m_{0} \alpha_{0}\left(i \cot \frac{\theta_{0}}{2}\right)^{\frac{2}{\pi} \alpha_{0}-1} \frac{-1}{2 \sin ^{2} \frac{\theta_{0}}{2}} .
\end{aligned}
$$

Using these equalities, we deduce

$$
f^{\prime}\left(z_{0}\right)=\left(i \cot \frac{\theta_{0}}{2}\right)^{\frac{2}{\pi} \alpha_{0}}\left(1+i \frac{2}{\pi} \alpha_{0} \frac{m_{0}}{\sin \theta_{0}}\right) .
$$

Thus, if $\theta_{0} \in[0, \pi]$, then

$$
\begin{equation*}
\left|\arg f^{\prime}\left(z_{0}\right)\right|=\alpha_{0}+\arctan \left(\frac{2}{\pi} \alpha_{0} \frac{m_{0}}{\sin \theta_{0}}\right) \geq \alpha_{0}+\arctan \left(\frac{2}{\pi} \alpha_{0}\right)>\frac{7 \pi}{10} \tag{12}
\end{equation*}
$$

and the case $\theta \in[-\pi, 0]$ is analogous to the previous one. If $\alpha_{0}=\frac{50 \pi}{108}$, then (12) holds, and this contradicts (10). The contradiction shows that inequality (11) holds.

We prove in the followings that

$$
\begin{equation*}
\left|\arg \frac{\mathrm{F}(z)}{z}\right|<\alpha_{1}=\frac{3 \pi}{10} z \in \mathrm{U} . \tag{13}
\end{equation*}
$$

This inequality is equivalent to the subordination

$$
p(z)=\frac{F(z)}{z} \prec\left(\frac{1+z}{1-z}\right)^{\frac{2}{\pi} \alpha_{1}}=q(z), z \in U .
$$

If this subordination does not hold, then we use again Lemma 1 and we get that there are two points $z_{1} \in U, \zeta_{1}=e^{i \theta_{1}} \in \partial U$ and a real number $m_{1} \in[1, \infty)$, such that

$$
\begin{aligned}
p\left(z_{1}\right) & =q\left(\zeta_{1}\right), \\
z_{1} p^{\prime}\left(z_{1}\right) & =m_{1} \zeta_{1} q^{\prime}\left(\zeta_{1}\right) .
\end{aligned}
$$

These equalities imply

$$
\begin{align*}
\frac{f\left(z_{1}\right)}{z_{1}} & =z_{1} p^{\prime}\left(z_{1}\right)+p\left(z_{1}\right) \\
& =\left(i \cot \frac{\theta_{1}}{2}\right)^{\frac{2}{\pi} \alpha_{1}}-\frac{2}{\pi} \alpha_{1} m_{1} \times\left(i \cot \frac{\theta_{1}}{2}\right)^{\frac{2}{\pi} \alpha_{1}-1} \frac{1}{2 \sin ^{2} \frac{\theta_{1}}{2}}  \tag{14}\\
& =\left(i \cot \frac{\theta_{1}}{2}\right)^{\frac{2}{\pi} \alpha_{1}}\left(1+i \frac{2}{\pi} \alpha_{1} \frac{m_{1}}{\sin \theta_{1}}\right)
\end{align*}
$$

If $\theta_{1} \in[0, \pi]$, then

$$
\arg \left(1+i \frac{2}{\pi} \alpha_{1} \frac{m_{1}}{\sin \theta_{1}}\right)=\arctan \left[\frac{2}{\pi} \frac{m_{1}}{\sin \theta_{1}}\right] \geq \arctan \left(\frac{2}{\pi} \alpha_{1}\right)
$$

and (14) implies

$$
\left|\arg \frac{f\left(z_{1}\right)}{z_{1}}\right| \geq \alpha_{1}+\arctan \left(\frac{2}{\pi} \alpha_{1}\right)>\frac{50 \pi}{108} .
$$

If $\theta \in[-\pi, 0]$, then the same inequality can be deduced. This contradicts (11) and the contradiction implies (13). Now we are able to prove that $F=A(f) \in$ $S^{*}$. Differentiating the equality $F=A(f)$ twice, we obtain

$$
\mathrm{F}^{\prime}(z)+z \mathrm{~F}^{\prime \prime}(z)=\mathrm{f}^{\prime}(z)
$$

This can be rewritten using the notations $\mathrm{p}(z)=\frac{z \mathrm{~F}^{\prime}(z)}{\mathrm{F}(z)}, \mathrm{P}(z)=\frac{\mathrm{F}(z)}{z g^{\prime}(z)}$ in the form

$$
\mathrm{P}(z)\left(z p^{\prime}(z)+p^{2}(z)\right)=\frac{f^{\prime}(z)}{g^{\prime}(z)} z \in U
$$

The conditions of the theorem imply that

$$
\begin{equation*}
\mathfrak{R}\left[\mathrm{P}(z)\left(z \mathrm{p}^{\prime}(z)+\mathrm{p}^{2}(z)\right)\right]>0, z \in \mathrm{U} \tag{15}
\end{equation*}
$$

We observe that (9) and (13) imply $|\arg (P(z))|<\frac{\pi}{2}, \quad z \in U$ and this is equivalent to $\mathfrak{R P}(z)>0, z \in U$. If $\mathfrak{R p}(z)>0, z \in U$ is not true, then according to Lemma 2 it follows that there are two real numbers $x_{2}, y_{2} \in \mathbb{R}$ and a point $z_{2} \in U$, such that $p\left(z_{2}\right)=\mathfrak{i x} x_{2}$ and $z_{2} p^{\prime}\left(z_{2}\right)=y_{2} \leq-\frac{1}{2}\left(x_{2}^{2}+1\right)$. Thus the equality

$$
\mathrm{P}\left(z_{2}\right)\left(z_{2} \mathrm{p}^{\prime}\left(z_{2}\right)+\mathrm{p}^{2}\left(z_{2}\right)\right)=\mathrm{P}\left(z_{2}\right)\left(\mathrm{y}_{2}-x_{2}^{2}\right)
$$

and $\mathfrak{R P}\left(z_{2}\right)>0$ imply that

$$
\mathfrak{R}\left[\mathrm{P}\left(z_{2}\right)\left(z_{2} \mathrm{p}^{\prime}\left(z_{2}\right)+\mathrm{p}^{2}\left(z_{2}\right)\right)\right] \leq 0 .
$$

This inequality contradicts (15), hence we deduce $\mathfrak{R} p(z)=\mathfrak{R} \frac{z \mathrm{~F}^{\prime}(z)}{\mathrm{F}(z)}>0, z \in \mathrm{U}$.

We end the paper stating a hypothesis.
Conjecture 1 We think that Theorem 3 and Theorem 4 can be improved in such a way that the obtained result would become an improvement of the second part of Theorem 1.

## References

[1] A. Imre, P. A. Kupán, R. Szász, Improvement of a criterion for starlikeness, Rocky Mountain J. Math., 42 (2012).
[2] S. S. Miller, P. T. Mocanu, Differential subordinations theory and applications, Marcel Dekker, New York, Basel, 2000.
[3] P. A. Kupán, R. Szász, About a condition for starlikeness, Ann. Univ. Sci. Budapest, Sect. Comp., 37 (2012), 261-274
[4] R. Szász, A counter-example concerning starlike functions, Stud. Univ. Babeş-Bolyai Math., LII. (2007), 171-172.
[5] R. Szász, An improvement of a criterion for starlikeness, Math. Pannon., 20/1 (2009), 69-77.

# A survey of the alternating sum-of-divisors function 

László Tóth<br>Universität für Bodenkultur, Institute of Mathematics<br>Department of Integrative Biology Gregor Mendel-Straße 33, A-1180 Wien, Austria and University of Pécs, Department of Mathematics<br>Ifjúság u. 6, H-7624 Pécs, Hungary<br>email: ltoth@gamma.ttk.pte.hu

Dedicated to the memory of my friend and colleague, Professor Antal Bege


#### Abstract

We survey arithmetic and asymptotic properties of the alternating sum-of-divisors function $\beta$ defined by $\beta\left(p^{a}\right)=p^{a}-p^{a-1}+$ $p^{a-2}-\cdots+(-1)^{a}$ for every prime power $p^{a}(a \geq 1)$, and extended by multiplicativity. Certain open problems are also stated.


## 1 Introduction

Let $\beta$ denote the multiplicative arithmetic function defined by $\beta(1)=1$ and

$$
\begin{equation*}
\beta\left(p^{a}\right)=p^{a}-p^{a-1}+p^{a-2}-\cdots+(-1)^{a} \tag{1}
\end{equation*}
$$

for every prime power $p^{a}(a \geq 1)$. That is,

$$
\begin{equation*}
\beta(n)=\sum_{d \mid n} d \lambda(n / d) \tag{2}
\end{equation*}
$$

for every integer $n \geq 1$, where $\lambda(n)=(-1)^{\Omega(n)}$ is the Liouville function, $\Omega(n)$ denoting the number of prime power divisors of $n$.

Key words and phrases: sum-of-divisors function, Liouville function, Euler's totient function, specially multiplicative function, imperfect number

The function $\beta$, as a variation of the sum-of-divisors function $\sigma$, was considered by Martin [19], Iannucci [16], Zhou and Zhu [32] regarding the following problem. In analogy with the perfect numbers, $\mathfrak{n}$ is said to be imperfect if $2 \beta(n)=n$. More generally, $n$ is said to be k-imperfect if $k \beta(n)=n$ for some integer $k \geq 2$. The only known imperfect numbers are $2,12,40,252,880$, 10880,75852 , 715816960 and 3074457344902430720 (sequence A127725 in [33]). Examples of 3-imperfect numbers are $6,120,126,2520$. No k-imperfect numbers are known for $k>3$. See also the book of Guy [9, p. 72].

This function occurs in the literature also in another context. Let

$$
\mathrm{b}(\mathrm{n})=\#\{\mathrm{k}: 1 \leq \mathrm{k} \leq \mathrm{n} \text { and } \operatorname{gcd}(\mathrm{k}, \mathrm{n}) \text { is a square }\}
$$

Then $b(n)=\beta(n)(n \geq 1)$, see Cohen [5, Cor. 4.2], Sivaramakrishnan [23], [24, p. 201], McCarthy [20, Sect. 6], [22, p. 25], Bege [3, p. 39], Iannucci [16, p. 12]. A modality to show the identity $b(n)=\beta(n)$ is to apply a familiar property of the Liouville function, namely,

$$
\begin{equation*}
\sum_{d \mid n} \lambda(d)=\chi(n) \quad(n \geq 1) \tag{3}
\end{equation*}
$$

where $\chi$ is the characteristic function of the set of squares. Using (3),

$$
\begin{aligned}
b(n)=\sum_{k=1}^{n} \chi(\operatorname{gcd}(k, n)) & =\sum_{k=1}^{n} \sum_{d \mid g c d(k, n)} \lambda(d)=\sum_{d \mid n} \lambda(d) \sum_{1 \leq k \leq n, d k} 1 \\
& =\sum_{d \mid n} \lambda(d) \frac{n}{d}=\beta(n) .
\end{aligned}
$$

In this paper we survey certain known properties of the function $\beta$, and give also other ones (without references), which may be known, but we could not locate them in the literature.

We point out that the function $\beta$ has a double character. On the one hand, certain properties of this function are similar to those of the sum-of-divisors function $\sigma$, due to the fact that both $\beta$ and $\sigma$ are the Dirichlet convolution of two completely multiplicative functions. Such functions are called in the literature specially multiplicative functions or quadratic functions. Their study in connection to the Busche-Ramanujan identities goes back to the work of Vaidyanathaswamy [30]. See also [15, 21, 22, 24].

On the other hand, further properties of this function are analogous to those of the Euler's totient function $\varphi$, as a consequence of the representation of $\beta$ given above.

We call the function $\beta$ the alternating sum-of-divisors function or alternating sigma function. Sivaramakrishnan [24, p. 201] remarked that it may be termed the square totient function.

It is possible, of course, to define other alternating sums of the positive divisors of $n$. For example, let $\theta(n)=\sum_{d \mid n} d \lambda(d)(n \geq 1)$. Then $\theta(n)=$ $\lambda(n) \beta(n)(n \geq 1)$. This is sequence A061020 in [33]. Another example: let $n=d_{1}>d_{2}>\cdots>d_{\tau(n)}=1$ be the divisors of $n$, in decreasing order, and let $A(n)=\sum_{j=1}^{\tau(n)}(-1)^{j-1} d_{j}$, cf. [2]. Note that the function $A$ is not multiplicative.

We do not give detailed proofs, excepting the proofs of formulae (10), (16), (17) and of the Proposition in Section 7, which are included in Section 8. We leave to the interested reader to compare the corresponding properties of the functions $\beta, \sigma$ and $\varphi$. See, for example, the books [1, 10, 22, 24, 27].

In Section 7 we pose certain open problems. One of them is concerning super-imperfect numbers $n$, defined by $2 \beta(\beta(n))=n$. This notion seems not to appear in the literature. The super-imperfect numbers up to $10^{7}$ are $\mathfrak{n}=$ $2,4,8,128,32768$. The number 2147483648 is also super-imperfect.

The corresponding concept for the sigma function is the following: A number n is called superperfect if $\sigma(\sigma(n))=2 n$. The even superperfect numbers are $2^{\mathrm{p}-1}$, where $2^{\mathrm{p}}-1$ is a Mersenne prime, cf. [26] (sequence A019279 in [33]). No odd superperfect numbers are known.

## 2 Basic properties

It is clear from (1) that for every prime power $p^{a}(a \geq 1)$,

$$
\beta\left(p^{a}\right)=\frac{p^{a+1}+(-1)^{a}}{p+1}= \begin{cases}\frac{p^{a+1}-1}{p+1}, & \text { if } a \geq 1 \text { is odd }  \tag{4}\\ \frac{p^{a+1}+1}{p+1}, & \text { if } a \geq 2 \text { is even }\end{cases}
$$

We obtain from (2),

$$
\sum_{n=1}^{\infty} \frac{\beta(n)}{n^{s}}=\frac{\zeta(s-1) \zeta(2 s)}{\zeta(s)} \quad(\Re(s)>2)
$$

where $\zeta$ is the Riemann zeta function, leading to another convolution representation of $\beta$, namely

$$
\begin{equation*}
\beta(n)=\sum_{d^{2} k=n} \varphi(k) \quad(n \geq 1) \tag{5}
\end{equation*}
$$

cf. McCarthy [20, Sect. 6], [22, p. 25], Bege [3, p. 39].

We have $\varphi(n) \leq \beta(n) \leq n(n \geq 1)$. More exactly, it follows from (5) that for every $n \geq 1$,

$$
\beta(n)=\varphi(n)+\sum_{d^{2} k=n, d>1} \varphi(k) \geq \varphi(n)
$$

with equality for the squarefree values of $n$. Also,

$$
\beta(n) \leq \sum_{d k=n} \varphi(k)=n
$$

with equality only for $n=1$.
Moreover, $\beta(n) \leq \varphi^{*}(n)$ for every $n \geq 1$, with equality if and only if n is squarefree or twice a squarefree number. This follows easily from (4). Here $\varphi^{*}$ is the unitary Euler function, which is multiplicative and given by $\varphi^{*}\left(p^{a}\right)=p^{a}-1$ for every prime power $p^{a}(a \geq 1)$, cf. [22, 24]. Also, $\beta(n) \geq \sqrt{n}$ $(n \geq 1, n \neq 2, n \neq 6)$.

Similar to the corresponding property of the function $\sigma, \beta(n)$ is odd if and only if $n$ is a square or twice a square.

The function $\beta$ appears in certain elementary identities regarding the set of squares, for example in

$$
\sum_{\substack{k=1 \\ \operatorname{gcd}(k, n) \text { a square }}}^{n} k=\frac{n(\beta(n)+\chi(n))}{2} \quad(n \geq 1)
$$

$$
\prod_{\substack{k=1 \\ n}} k=n^{\beta(n)} \prod_{d \mid n}\left(d!/ d^{d}\right)^{\lambda(n / d)} \quad(n \geq 1)
$$

$$
\operatorname{gcd}(k, n) \text { a square }
$$

which can be deduced from (3).

## 3 Generalizations

An obvious generalization of $\beta$ is the function $\beta_{a}(a \in \mathbb{C})$ defined by

$$
\begin{equation*}
\beta_{a}(n)=\sum_{d \mid n} d^{a} \lambda(n / d) \quad(n \geq 1) \tag{6}
\end{equation*}
$$

If $a=m$ is a positive integer, then the following representation can be given: $\beta_{m}(n)=\#\left\{k: 1 \leq k \leq n^{m},\left(k, n^{m}\right)_{m}\right.$ is a $2 m$-th power $\}$, where $(a, b)_{m}$
stands for the largest common $m$-th power divisor of $a$ and $b$. See McCarthy [20, Sect. 6], [22, p. 51].

Note that if $a=0$, then $\beta_{0}=\chi$, the characteristic function of the set of squares, used above.

For an arbitrary nonempty set $S$ of positive integers let $\varphi_{S}(n)=\#\{k: 1 \leq$ $k \leq n, \operatorname{gcd}(k, n) \in S\}$. For $S=\{1\}$ and $S$ the set of squares this reduces to Euler's function $\varphi$ and to the function $\beta$, respectively. The function $\varphi_{S}$ was investigated by Cohen [6]. For every set $S$ one has

$$
\varphi_{S}(n)=\sum_{d \mid n} d \mu_{S}(n / d) \quad(n \geq 1)
$$

where the function $\mu_{S}$ is defined by $\sum_{d \mid n} \mu_{S}(d)=\chi_{S}(n)(n \geq 1)$, i.e., $\mu_{S}=\mu_{*}$ $\chi_{\mathrm{s}}$ in terms of the Dirichlet convolution $*, \chi_{\mathrm{s}}$ and $\mu$ denoting the characteristic function of $S$ and the Möbius function, respectively.

Also, let

$$
\mathrm{B}(\mathrm{r}, \mathfrak{n})=\sum_{\substack{k=1 \\ \operatorname{gcd}(k, n) \text { a square }}}^{n} \exp (2 \pi i k r / \mathfrak{n}),
$$

which is an analog of the Ramanujan sum to be considered in Section 5. Then

$$
B(r, n)=\sum_{d \mid g c d(r, n)} d \lambda(n / d) \quad(r, n \geq 1)
$$

see Sivaramakrishnan [23], [24, p. 202], Haukkanen [13]. For $r=n$ one has $B(n, n)=\beta(n)$.
These generalizations can also be combined. See also Haukkanen [11, 12]. Many of the results given in the present paper can be extended for these generalizations.

We consider in what follows only the functions $\beta_{\mathrm{a}}$ defined by (6) and do not deal with other generalizations.

## 4 Further properties

For every $n, m \geq 1$,

$$
\begin{equation*}
\beta(n) \beta(m)=\sum_{d \mid \operatorname{cd}(n, m)} \beta\left(n m / d^{2}\right) d \lambda(d) \tag{7}
\end{equation*}
$$

and equivalently,

$$
\begin{equation*}
\beta(n m)=\sum_{d \mid \operatorname{gcd}(n, m)} \beta(n / d) \beta(m / d) d \mu^{2}(d) \tag{8}
\end{equation*}
$$

cf. [23], [22, p. 26]. Here (7) and (8) are special cases of the Busche-Ramanujan identities, valid for specially multiplicative functions. See [15, 21, 22, 24, 30] for their discussions and proofs.

Direct proofs of (7) and (8) can be given by showing that both sides of these identities are multiplicative, viewed as functions of two variables and then computing their values for prime powers. Recall that an arithmetic function $f$ of two variables is called multiplicative if it is nonzero and $f\left(n_{1} m_{1}, n_{2} m_{2}\right)=$ $f\left(n_{1}, n_{2}\right) f\left(m_{1}, m_{2}\right)$ holds for any $n_{1}, n_{2}, m_{1}, m_{2} \geq 1$ such that $\operatorname{gcd}\left(n_{1} n_{2}\right.$, $\left.m_{1} m_{2}\right)=1$. See [30], [29], [24, Ch. VII].

The proof of the equivalence of identities of type (7) and (8) is outlined in [14], referring to the work of Vaidyanathaswamy [30].

It follows at once from (8) that $\beta(n m) \geq \beta(n) \beta(m)$ for every $n, m \geq 1$, i.e., $\beta$ is super-multiplicative. Formula (8) leads also to the double Dirichlet series

$$
\sum_{n, m=1}^{\infty} \frac{\beta(n m)}{n^{s} m^{t}}=\frac{\zeta(s-1) \zeta(2 s) \zeta(t-1) \zeta(2 t) \zeta(s+t-1)}{\zeta(s) \zeta(t) \zeta(2(s+t-1))}
$$

valid for $s, t \in \mathbb{C}$ with $\mathfrak{R}(s)>2, \mathfrak{R}(\mathrm{t})>2$.
The generating power series of $\beta$ is

$$
\sum_{n=1}^{\infty} \beta(n) x^{n}=\sum_{n=1}^{\infty} \frac{\lambda(n) x^{n}}{\left(1-x^{n}\right)^{2}} \quad(|x|<1)
$$

which is a direct consequence of (2).
Consider the functions $\beta_{a}$ defined by (6). One has

$$
\begin{equation*}
\sum_{n=1}^{\infty} \frac{\beta_{a}(n) \beta_{b}(n)}{n^{s}}=\frac{\zeta(s) \zeta(s-a-b) \zeta(2 s-2 a) \zeta(2 s-2 b)}{\zeta(s-a) \zeta(s-b) \zeta(2 s-a-b)} \tag{9}
\end{equation*}
$$

valid for every $s, a, b \in \mathbb{C}$ with $\mathfrak{R}(s)>1+\max (0, \mathfrak{R}(a), \mathfrak{R}(b), \mathfrak{R}(a+b))$.
This formula is similar to Ramanujan's well-known result for the product $\sigma_{a}(n) \sigma_{b}(n)$, where $\sigma_{a}(n)=\sum_{d \mid n} d^{a}$. Formula (9) is due to Chowla [4], in an equivalent form for the product $\theta_{a}(n) \theta_{b}(n)$, where $\theta_{a}(n)=\sum_{d \mid n} d^{a} \lambda(d)$.

Formula (9) and that of Ramanujan follow from the next more general result concerning the product of two arbitrary specially multiplicative functions.

If $f, g, h, k$ are completely multiplicative functions, then

$$
\begin{equation*}
(f * g)(h * k)=f h * f k * g h * g k * w \tag{10}
\end{equation*}
$$

where $w(n)=\mu(m) f(m) g(m) h(m) k(m)$ if $n=m^{2}$ is a square and $w(n)=0$ otherwise.

This result is given by Vaidyanathaswamy [30, p. 621], Lambek [17], Subbarao [25]. See also [24, p. 50]. The proof of (10) can be carried out using Euler products. This is well-known in the case of Ramanujan's result regarding $\sigma_{\mathrm{a}} \sigma_{\mathrm{b}}$, and is presented in several texts, cf., e.g., [10, Th. 305], [22, Prop. 5.4]. An alternative proof is given by Lambek [17].

In Section 8 we offer another less known short proof of (10).
In the case of the functions $f(n)=n^{a}, h(n)=n^{b}, g(n)=k(n)=\lambda(n)$ we obtain (9) by using the known formulae for the Dirichlet series corresponding to the right hand side of (10).

If $f(n)=n^{a}, h(n)=n^{b}, g(n)=\lambda(n), k(n)=1$, then we deduce

$$
\sum_{n=1}^{\infty} \frac{\beta_{a}(n) \sigma_{b}(n)}{n^{s}}=\frac{\zeta(s-a) \zeta(s-a-b) \zeta(2 s) \zeta(2 s-2 b) \zeta(2 s-a-b)}{\zeta(s) \zeta(s-b) \zeta(4 s-2 a-2 b)},
$$

valid for the same region as (9).
Remark that we obtain, as direct corollaries, the next formulae:

$$
\begin{gather*}
\sum_{n=1}^{\infty} \frac{\beta^{2}(n)}{n^{s}}=\frac{\zeta(s) \zeta(s-2) \zeta(2 s-2)}{\zeta^{2}(s-1)}, \\
\sum_{n=1}^{\infty} \frac{\beta\left(n^{2}\right)}{n^{s}}=\frac{\zeta(s) \zeta(s-2)}{\zeta(s-1)},  \tag{11}\\
\sum_{n=1}^{\infty} \frac{\beta(n) \sigma(n)}{n^{s}}=\frac{\zeta(s-2) \zeta(2 s) \zeta^{2}(2 s-2)}{\zeta(s) \zeta(4 s-4)},
\end{gather*}
$$

all valid for $\mathfrak{R}(s)>3$. Here (11) is obtained from (9) by choosing $a=1$ and $\mathrm{b}=0$.

From these Dirichlet series representations we can deduce the following convolutional identities:

$$
\begin{gather*}
\beta^{2}(n)=\sum_{d k=n} d 2^{\omega(d)} \lambda(d) \sigma_{2}(k) \quad(n \geq 1),  \tag{12}\\
\beta\left(n^{2}\right)=\sum_{d k=n} d \mu(d) \sigma_{2}(k) \quad(n \geq 1), \tag{13}
\end{gather*}
$$

$$
\begin{equation*}
\beta(n) \sigma(n)=\sum_{d^{2} k=n} d^{2} 2^{\omega(d)} \beta_{2}(k) \quad(n \geq 1) \tag{14}
\end{equation*}
$$

where $\omega(n)$ denotes the number of distinct prime factors of $n$.

## 5 Asymptotic behavior

The average order of $\beta(n)$ is $\left(\pi^{2} / 15\right) n$, more exactly,

$$
\begin{equation*}
\sum_{n \leq x} \beta(n)=\frac{\pi^{2}}{30} x^{2}+\mathcal{O}\left(x(\log x)^{2 / 3}(\log \log x)^{4 / 3}\right) \tag{15}
\end{equation*}
$$

Formula (15) follows from the convolution representation (5) and from the known estimate of Walfisz regarding $\sum_{n \leq x} \varphi(n)$ with the same error term as above.

There are also other asymptotic properties of the $\varphi$ function, which can be transposed to $\beta$ by using that $\beta(n) \geq \varphi(n)$, with equality for $n$ squarefree. For example,

$$
\liminf _{n \rightarrow \infty} \frac{\beta(n) \log \log n}{n}=e^{-\gamma}
$$

where $\gamma$ is Euler's constant (cf. [10, Th. 328] concerning $\varphi$ ). Another example: the set $\{\beta(n) / n: n \geq 1\}$ is dense in the interval $[0,1]$.

Let $c_{r}(n)$ denote the Ramanujan sum, defined as the sum of $n$-th powers of the primitive $r$-th roots of unity. Then

$$
\begin{gather*}
\frac{\beta(n)}{n}=\frac{\pi^{2}}{15} \sum_{r=1}^{\infty} \frac{\lambda(r)}{r^{2}} c_{r}(n)  \tag{16}\\
=\frac{\pi^{2}}{15}\left(1-\frac{(-1)^{n}}{2^{2}}-\frac{2 \cos (2 \pi n / 3)}{3^{2}}+\frac{2 \cos (\pi n / 2)}{4^{2}}+\ldots\right),
\end{gather*}
$$

showing how the values of $\beta(n) / n$ fluctuate harmonically about their mean value $\pi^{2} / 15$, cf. [7], [22, p. 245].

A quick direct proof of formula (16) is given in Section 8. We refer to [18] for a recent survey of expansions of functions with respect to Ramanujan sums.

From the identities (12), (13) and (14) we deduce the following asymptotic formulae:

$$
\sum_{n \leq x} \beta^{2}(n)=\frac{2 \zeta(3)}{15} x^{3}+\mathcal{O}\left(x^{2}(\log x)^{2}\right)
$$

$$
\begin{aligned}
& \sum_{n \leq x} \beta\left(n^{2}\right)=\frac{2 \zeta(3)}{\pi^{2}} x^{3}+\mathcal{O}\left(x^{2} \log x\right), \\
& \sum_{n \leq x} \beta(n) \sigma(n)=\frac{\pi^{6}}{2430 \zeta(3)} x^{3}+\mathcal{O}\left(x^{2}\right) .
\end{aligned}
$$

We also have

$$
\begin{equation*}
\sum_{n \leq x} \frac{1}{\beta(n)}=K_{1} \log x+K_{2}+\mathcal{O}\left(x^{-1+\varepsilon}\right) \tag{17}
\end{equation*}
$$

for every $\varepsilon>0$, where $K_{1}$ and $K_{2}$ are constants,

$$
K_{1}=\prod_{p}\left(1-\frac{1}{p}\right)\left(1+\sum_{a=1}^{\infty} \frac{1}{\beta\left(p^{a}\right)}\right) .
$$

For the proof of (17) see Section 8.

## 6 Unitary analog

Consider the function $\beta^{*}$ defined by

$$
\beta^{*}(n)=\sum_{d \| n} d \lambda(n / d) \quad(n \geq 1)
$$

where the sum is over the unitary divisors $d$ of $n$. Recall that $d$ is a unitary divisor of $n$ if $d \mid n$ and $\operatorname{gcd}(d, n / d)=1$. Here $\beta^{*}\left(p^{a}\right)=p^{a}+(-1)^{a}$ for every prime power $p^{a}(a \geq 1)$ and

$$
\begin{equation*}
\sum_{n=1}^{\infty} \frac{\beta^{*}(n)}{n^{s}}=\frac{\zeta(s-1) \zeta(2 s) \zeta(2 s-1)}{\zeta(s) \zeta(4 s-2)} \quad(\mathfrak{R}(s)>2) . \tag{18}
\end{equation*}
$$

The formula (18) can be derived using Euler products or by establishing the convolutional identity

$$
\beta^{*}(n)=\sum_{d k^{2}=n} \beta(d) k q(k),
$$

q standing for the characteristic function of the squarefree numbers. This leads also to the asymptotic formula

$$
\sum_{n \leq x} \beta^{*}(n)=\frac{63 \zeta(3)}{2 \pi^{4}} x^{2}+\mathcal{O}\left(x(\log x)^{5 / 3}(\log \log x)^{4 / 3}\right) .
$$

Note the following interpretation: $\beta^{*}(n)=\#\left\{k: 1 \leq k \leq n\right.$ and $\operatorname{gcd}(k, n)_{*}$ is a square\}, where $(a, b)_{*}$ is the largest divisor of $a$ which is a unitary divisor of $b$.

## 7 Super-imperfect numbers, open problems

A number $n$ is super-imperfect if $2 \beta(\beta(n))=n$, cf. Introduction. Observe that, excepting 4 , all the other examples of super-imperfect numbers are of the form $\mathrm{n}=2^{2^{k}-1}$ with $\mathrm{k} \in\{1,2,3,4,5\}$. The proof of the next statement is given in Section 8.

Proposition. For $\mathrm{k} \geq 1$ the number $\mathrm{n}_{\mathrm{k}}=2^{2^{k}-1}$ is super-imperfect if and only if $k \in\{1,2,3,4,5\}$.

Problem 1. Is there any other super-imperfect number?
More generally, we define $n$ to be (m,k)-imperfect if $k \beta^{(m)}(n)=n$, where $\beta^{(m)}$ is the $m$-fold iterate of $\beta$. For example, $3,15,18,36,72,255$ are (2, 3)imperfect, 6, 12, 24, 30, 60, 120 are (2,6)-imperfect, 6, 36, 144 are (3,6)imperfect numbers.

We refer to [8] regarding ( $m, k$ )-perfect numbers, defined by $\sigma^{(m)}(n)=k n$.
Problem 2. Investigate the ( $\mathrm{m}, \mathrm{k}$ )-imperfect numbers.
The numbers $n=1,20,116,135,171,194,740, \ldots$ are solutions of the equation $\beta(n)=\beta(n+1)$.

Problem 3. Are there infinitely many numbers $n$ such that $\beta(n)=\beta(n+1)$ ?
Remark that it is not known if there are infinitely many numbers $n$ such that $\sigma(n)=\sigma(n+1)$ (sequence A002961 in [33]). See also Weingartner [31].

The next problem is the analog of Lehmer's open problem concerning the $\varphi$ function.

Problem 4. Is there any composite number $n \neq 4$ such that $\beta(n)$ divides $\mathrm{n}-1$ ?

Up to $10^{6}$ there are no such composite numbers.
The computations were performed using Maple. The function $\beta(n)$ was generated by the following procedure:

```
beta:= proc(n) local x, i: x:= 1:
for i from 1 to nops(ifactors(n)[2])
do p_i:= ifactors(n)[2][i][1]: a_i:= ifactors(n)[2][i] [2];
x:= x*((p_i^(a_i+1)+(-1)^(a_i))/(p_i+1)): od: RETURN(x) end;
# alternating sum-of-divisors function
```


## 8 Proofs

Proof of formula (10): Write

$$
(f * g)(n)(h * k)(n)=\sum_{\substack{d|n \\ e| n}} f(d) g(n / d) h(e) k(n / e)
$$

where $d|n, e| n \Leftrightarrow \operatorname{lcm}[d, e] \mid n$. Write $d=m u, e=m v$ with $\operatorname{gcd}(u, v)=1$. Then $\operatorname{lcm}[\mathrm{d}, \mathrm{e}]=\mathfrak{m u v}$ and obtain that this sum is

$$
\begin{aligned}
& \quad \sum_{\substack{m u v \mathfrak{n} \\
\operatorname{gcd}(u, v)=1}} f(m u) g(n /(m u)) h(m v) k(n /(m v)) \\
& =\sum_{m u \cup n} f(m u) g(n /(m u)) h(m v) k(n /(m v)) \sum_{\delta \mid(\operatorname{gcd}(u, v)} \mu(\delta) .
\end{aligned}
$$

Putting now $u=\delta x, v=\delta y$ and using that the considered functions are all completely multiplicative the latter sum is

$$
\sum_{\delta^{2} x y m=n}(\mu f g h k)(\delta)(f k)(x)(g h)(y)(f h)(m)(g k)(t)
$$

finishing the proof (cf. [30, p. 621] and [27, p. 161]).
Proof of formula (16): Let $\eta_{\mathrm{r}}(\mathrm{n})=\mathrm{r}$ if $\mathrm{r} \mid \mathrm{n}$ and $\eta_{\mathrm{r}}(\mathrm{n})=0$ otherwise. Applying that $\sum_{d \mid r} c_{d}(n)=\eta_{r}(n)$ we deduce

$$
\begin{aligned}
\frac{\beta(n)}{n} & =\sum_{d \mid n} \frac{\lambda(d)}{d}=\sum_{d=1}^{\infty} \frac{\lambda(d)}{d^{2}} \eta_{d}(n)=\sum_{d=1}^{\infty} \frac{\lambda(d)}{d^{2}} \sum_{r \mid d} c_{r}(n) \\
& =\sum_{r=1}^{\infty} \frac{\lambda(r)}{r^{2}} c_{r}(n) \sum_{k=1}^{\infty} \frac{\lambda(k)}{k^{2}}=\frac{\zeta(4)}{\zeta(2)} \sum_{r=1}^{\infty} \frac{\lambda(r)}{r^{2}} c_{r}(n)
\end{aligned}
$$

using that $\lambda$ is completely multiplicative and its Dirichlet series is $\sum_{n=1}^{\infty} \lambda(n) / n^{s}$ $=\zeta(2 s) / \zeta(s)$. The rearranging of the terms is justified by the absolute convergence.

Proof of formula (17): Write

$$
\frac{1}{\beta(n)}=\sum_{\substack{d k=n \\ \operatorname{gcd}(d, k)=1}} \frac{h(d)}{\varphi(k)}
$$

as the unitary convolution of the functions $h$ and $1 / \varphi$, where $h$ is multiplicative and for every prime power $p^{a}(a \geq 1)$,

$$
\frac{1}{\beta\left(p^{a}\right)}=h\left(p^{a}\right)+\frac{1}{\varphi\left(p^{a}\right)}, \quad h\left(p^{a}\right)=-\frac{p^{a-1}+(-1)^{a}}{p^{a-1}(p-1)\left(p^{a+1}+(-1)^{a}\right)} .
$$

Here

$$
\left|h\left(p^{a}\right)\right| \leq \frac{1}{p^{a}(p-1)^{2}}, \quad|h(n)| \leq \frac{f(n)}{\varphi(n)}(n \geq 1)
$$

with $f(n)=\prod_{p \mid n}(p(p-1))^{-1}$. We deduce

$$
\sum_{n \leq x} \frac{1}{\beta(n)}=\sum_{d \leq x} h(d) \sum_{\substack{k \leq x / d \\ \operatorname{gcd}(d, k)=1}} \frac{1}{\varphi(k)}
$$

and use the known estimates for the inner sum. The same arguments were applied in the proof of [28, Th. 2].

Proof of the Proposition of Section 7: The fact that the numbers $n_{k}$ with $1 \leq \mathrm{k} \leq 5$ are super-imperfect follow also by direct computations, but the following arguments reveal a connection to the Fermat numbers $F_{m}=2^{2^{m}}+1$.

For $n_{k}=2^{2^{k}-1}$ with $k \geq 1$ we have

$$
\beta\left(n_{k}\right)=\frac{2^{2^{k}}-1}{3}=F_{1} F_{2} \cdots F_{k-1}
$$

(for $k=1$ this is 1 , the empty product). Since the numbers $F_{m}$ are pairwise relatively prime,

$$
\beta\left(\beta\left(n_{k}\right)\right)=\beta\left(F_{1}\right) \beta\left(F_{2}\right) \cdots \beta\left(F_{k-1}\right) .
$$

Now for $2 \leq k \leq 5$, using that $F_{1}, F_{2}, F_{3}, F_{4}$ are primes,

$$
\beta\left(\beta\left(n_{k}\right)\right)=2^{2^{1}} \cdot 2^{2^{2}} \cdot \ldots \cdot 2^{2^{k-1}}=2^{2^{k}-2}=\frac{n_{k}}{2},
$$

showing that $n_{k}$ is super-imperfect.
Now let $k \geq 6$. We use that $F_{5}$ is composite and that $\beta(n) \supsetneqq n-1$ for every $n \neq 4$ composite. Hence $\beta\left(F_{5}\right) \nsupseteq 2^{2^{5}}$ and

$$
\beta\left(\beta\left(n_{k}\right)\right) \nLeftarrow \beta\left(F_{1}\right) \beta\left(F_{2}\right) \cdots \beta\left(F_{k-1}\right)=2^{2^{k}-2}=\frac{n_{k}}{2},
$$

ending the proof.
Note that for $k \geq 2$ the number $m_{k}=2^{2^{k}-1} F_{1} F_{2} \cdots F_{k-1}$ is imperfect if and only if $k \in\{2,3,4,5\}$. This follows by similar arguments. The imperfect numbers of this form are $40,10880,715816960$ and 3074457344902430720 , given in the Introduction.

## Acknowledgement

The author gratefully acknowledges support from the Austrian Science Fund (FWF) under the projects Nr. P20847-N18 and M1376-N18.

## References

[1] T. M. Apostol, Introduction to Analytic Number Theory, Springer, New York, Heidelberg, Berlin, 1976.
[2] K. Atanassov, A remark on an arithmetic function. Part 2, Notes Number Theory Discrete Math. (NNTDM), 15 (2009), no 3, 21-22.
[3] A. Bege, Old and New Arithmetic Functions (Hungarian), Scientia, Kolozsvár [Cluj-Napoca], 2006, 115 pp.
[4] S. D. Chowla, On some identities involving zeta-functions, J. Indian Math. Soc., 17 (1928), 153-163.
[5] E. Cohen, Representations of even functions (mod r), I. Arithmetical identities, Duke Math. J., 25 (1958), 401-421.
[6] E. Cohen, Arithmetical functions associated with arbitrary sets of integers, Acta Arith., 5 (1959), 407-415.
[7] E. Cohen, Fourier expansions of arithmetical functions, Bull. Amer. Math. Soc., 67 (1961), 145-147.
[8] G. L. Cohen, H. J. J. te Riele, Iterating the sum-of-divisors function, Exp. Math., 5 (1996), 91-100.
[9] R. K. Guy, Unsolved Problems in Number Theory, 3rd ed., SpringerVerlag, New York, Berlin, Heidelberg, London, Paris, Tokyo, 2004.
[10] G. H. Hardy, I. M. Wright, An Introduction to the Theory of Numbers, Oxford Univ. Press, Fourth ed., 1975.
[11] P. Haukkanen, Some generalized totient functions, Math. Stud., 56 (1989), 65-74.
[12] P. Haukkanen, An abstract Möbius inversion formula with numbertheoretic applications, Discrete Math., 142 (1995), 87-96.
[13] P. Haukkanen, On some set reduced arithmetical sums, Indian J. Math., 39 (1997), 147-158.
[14] P. Haukkanen, An exponential Busche-Ramanujan identity, Mathematica (Cluj), 41 (64) (1999), 177-185.
[15] P. Haukkanen, Some characterizations of specially multiplicative functions, Int. J. Math. Math. Sci., 37 (2003), 2335-2344.
[16] D. E. Iannucci, On a variation of perfect numbers, Integers: Electronic J. of Combinatorial Number Theory, 6 (2006), \#A41, 13 pp.
[17] J. Lambek, Arithmetical functions and distributivity, Amer. Math. Monthly, 73 (1966), 969-973.
[18] L. G. Lucht, A survey of Ramanujan expansions, Int. J. Number Theory, 6 (2010), 1785-1799.
[19] G. Martin, Problem 99:08, Western Number Theory Problems, Compiled by G. Myerson, 2000, http://www.math.colostate.edu/ achter/wntc/problems/problems2000.pdf
[20] P. J. McCarthy, Some remarks on arithmetical identities, Amer. Math. Monthly, 67 (1960), 539-548.
[21] P. J. McCarthy, Busche-Ramanujan identities, Amer. Math. Monthly, 67 (1960), 966-970.
[22] P. J. McCarthy, Introduction to Arithmetical Functions, Springer, 1986.
[23] R. Sivaramakrishnan, Square-reduced residue systems (mod r) and related arithmetical functions, Canad. Math. Bull., 22 (1979), 207-220.
[24] R. Sivaramakrishnan, Classical Theory of Arithmetic Functions, Marcel Dekker, New York, 1989.
[25] M. V. Subbarao, Arithmetic functions and distributivity, Amer. Math. Monthly, 75 (1968), 984-988.
[26] D. Suryanarayana, Super perfect numbers, Elem. Math., 24 (1969), 1617.
[27] G. Tenenbaum, Introduction to Analytic and Probabilistic Number Theory, Cambridge Univ. Press, 1995.
[28] L. Tóth, A survey of gcd-sum functions, J. Integer Sequences, 13 (2010), Article 10.8.1, 23 pp.
[29] L. Tóth, Menon's identity and arithmetical sums representing functions of several variables, Rend. Sem. Mat. Univ. Politec. Torino, 69 (2011), 97-110.
[30] R. Vaidyanathaswamy, The theory of multiplicative arithmetic functions, Trans. Amer. Math. Soc., 33 (1931), 579-662.
[31] A. Weingartner, On the solutions of $\sigma(\mathfrak{n})=\sigma(\mathfrak{n}+\mathrm{k})$, J. Integer Sequences, 14 (2011), Article 11.5.5, 7 pp.
[32] W. Zhou, L. Zhu, On k-imperfect numbers, Integers: Electronic J. of Combinatorial Number Theory, 9 (2009), \#A01, 8 pp.
[33] The On-Line Encyclopedia of Integer Sequences, http://oeis.org.

Received: 13 October 2012

## Acta Universitatis Sapientiae

The scientific journal of the Sapientia University publishes original papers and deep surveys in several areas of sciences written in English.
Information about the appropriate series can be found at the Internet address
http://www.acta.sapientia.ro.
Editor-in-Chief
László DÁVID
Main Editorial Board
Zoltán A. BIRÓ Zoltán KÁSA András KELEMEN
Ágnes PETHŐ Emőd VERESS

# Acta Universitatis Sapientiae, Mathematica 

Executive Editor<br>Róbert SZÁSZ (Sapientia University, Romania)<br>rszasz@ms.sapientia.ro<br>Editorial Board<br>Sébastien FERENCZI (Institut de Mathématiques de Luminy, France)<br>Kálmán GYŐRY (University of Debrecen, Hungary)<br>Zoltán MAKÓ (Sapientia University, Romania)<br>Ladislav MIŠÍK (University of Ostrava, Czech Republic) János TÓTH (Selye University, Slovakia)<br>Adrian PETRUŞEL (Babeş-Bolyai University, Romania)<br>Alexandru HORVÁTH (Petru Maior University of Tg. Mureş, Romania)<br>Árpád BARICZ (Babeş-Bolyai University, Romania)<br>Csaba SZÁNTÓ (Babeş-Bolyai University, Romania)<br>Assistant Editor<br>Pál KUPÁN (Sapientia University, Romania)<br>Contact address and subscription:<br>Acta Universitatis Sapientiae, Mathematica<br>RO 400112 Cluj-Napoca<br>Str. Matei Corvin nr. 4.<br>Email: acta-math@acta.sapientia.ro<br>Each volume contains two issues.<br><br>Sapientia University<br><br>Scientia Publishing House

ISSN 1844-6094
http://www.acta.sapientia.ro

## Information for authors

Acta Universitatis Sapientiae, Mathematica publishes original papers and surveys in all field of Mathematics. All papers will be peer reviewed.

Papers published in current and previous volumes can be found in Portable Document Format (pdf) form at the address: http://www.acta.sapientia.ro.

The submitted papers should not be considered for publication by other journals. The corresponding author is responsible for obtaining the permission of coauthors and of the authorities of institutes, if needed, for publication, the Editorial Board disclaims any responsibility.

Submission must be made by email (acta-math@acta.sapientia.ro) only, using the LaTeX style and sample file at the address: http://www.acta.sapientia.ro. Beside the LaTeX source a pdf format of the paper is needed too.

Prepare your paper carefully, including keywords, 2010 Mathematics Subject Classification (MSC 2010) codes (http://www.ams.org/msc//msc2010.html), and the reference citation should be written in brackets in the text as [3]. References should be listed alphabetically using the following examples:

For papers in journals:
A. Hajnal, V. T. Sós, Paul Erdős is seventy, J. Graph Theory, 7 (1983), 391-393. For books:
D. Stanton, D. White, Constructive combinatorics, Springer, New York, 1986. For papers in contributed volumes:
Z. Csörnyei, Compilers in Algorithms of informatics, Vol. 1. Foundations (ed. A. Iványi), mondAt Kiadó, Budapest, 2007, pp. 80-130.
For internet sources:
E. Ferrand, An analogue of the Thue-Morse sequence, Electron. J. Comb., 14 (2007) \#R30, http://www.combinatorics.org/.

Illustrations should be given in Encapsulated Postcript (eps) format.
Authors are encouraged to submit papers not exceeding 15 pages, but no more than 10 pages are preferable.

One issue is offered each author. No reprints will be available.

## Printed by Gloria Printing House

Director: Péter Nagy


[^0]:    2010 Mathematics Subject Classification: 51M16

[^1]:    2010 Mathematics Subject Classification: 16S40, 16S35, 18E05

