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# Journal of Mathematics and Applications 

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#### Abstract

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# On a starlikeness of order $\alpha$ condition for meromorphic $m$-valent functions 

Adriana Cătaş

Submitted by: Jan Stankiewicz

Abstract: The aim of the paper is to provide sufficient conditions for starlikeness of order $\alpha$ for meromorphic $m$-valent functions in the punctured disc. The present work is based on some results involving differential subordinations.

AMS Subject Classification: 30C45
Key Words and Phrases: Analytic functions, starlike functions, meromorphic mvalent functions

## 1 Introduction and preliminaries

Let $\Sigma_{m}$ denote the class of functions of the form

$$
\begin{equation*}
f(z)=\frac{1}{z^{m}}+\sum_{n=1}^{\infty} a_{m+n-1} z^{m+n-1}, \quad m \in \mathbb{N}^{*} \tag{1.1}
\end{equation*}
$$

which are analytic and $m$-valent in the punctured disc

$$
\dot{U}=\{z \in C: 0<|z|<1\}=U \backslash\{0\} .
$$

A function $f \in \Sigma_{m}$ is said [1]. to be in the class $\Omega_{m}(\alpha)$ of meromorphic $m$-valently starlike functions of order $\alpha$ in $\dot{U}$ if and only if

$$
\begin{equation*}
\operatorname{Re}\left\{-\frac{z f^{\prime}(z)}{f(z)}\right\}>\alpha, \quad z \in \dot{U}, 0 \leq \alpha<m, m \in \mathbb{N}^{*} \tag{1.2}
\end{equation*}
$$

We denote $\Omega_{m}(0)=\Omega_{m}^{*}$.
The following definitions and lemmas will be used in the next section.
Let $\mathcal{H}(U)$ denote the space of analytic functions in $U$. For $n$ a positive integer and $a \in \mathbb{C}$ let

$$
\begin{equation*}
\mathcal{H}_{n}=\left\{f \in \mathcal{H}(U): \quad f(z)=a_{n} z^{n}+a_{n+1} z^{n+1}+\ldots\right\} \tag{1.3}
\end{equation*}
$$

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and

$$
\begin{equation*}
\mathcal{H}[a, n]=\left\{f \in \mathcal{H}(U): \quad f(z)=a+a_{n} z^{n}+a_{n+1} z^{n+1}+\ldots\right\} . \tag{1.4}
\end{equation*}
$$

For two functions $f$ and $g$ analytic in $U$, we say that the function $f(z)$ is subordinate to $g(z)$ in $U$ and write

$$
f \prec g \quad \text { or } \quad f(z) \prec g(z), \quad z \in U
$$

if there exists a Schwarz function $w(z)$, analytic in $U$ with

$$
w(0)=0 \quad \text { and } \quad|w(z)|<1, \quad z \in U
$$

such that

$$
\begin{equation*}
f(z)=g(w(z)), \quad z \in U \tag{1.5}
\end{equation*}
$$

In particular, if the function $g$ is univalent in $U$, the above subordination is equivalent to

$$
f(0)=g(0) \quad \text { and } \quad f(U) \subset g(U)
$$

Lemma 1.1 [2] Let $m$ be a positive integer and let $\alpha$ be real, with $0 \leq \alpha<m$. Let $q \in \mathcal{H}(U)$, with $q(0)=0, q^{\prime}(0) \neq 0$ and

$$
\begin{equation*}
\operatorname{Re}\left(1+\frac{z q^{\prime \prime}(z)}{q^{\prime}(z)}\right)>\frac{\alpha}{m} \tag{1.6}
\end{equation*}
$$

Define the function $h$ as

$$
\begin{equation*}
h(z)=m z q^{\prime}(z)-\alpha q(z) \tag{1.7}
\end{equation*}
$$

If $p \in \mathcal{H}_{m}$ and

$$
\begin{equation*}
z p^{\prime}(z)-\alpha p(z) \prec h(z) \tag{1.8}
\end{equation*}
$$

then $p(z) \prec q(z)$ and this result is sharp.
Lemma $1.2[3]$ Let $n \in \mathbb{N}^{*}$, let $\alpha \in[0,1]$ and let

$$
\begin{equation*}
M_{n}(\alpha)=\frac{n+1-\alpha}{\sqrt{(n+1-\alpha)^{2}+\alpha^{2}}+1-\alpha} \tag{1.9}
\end{equation*}
$$

If the function $f(z)$ of the form

$$
\begin{equation*}
f(z)=\frac{1}{z}+\sum_{k=n}^{\infty} a_{k} z^{k} \tag{1.10}
\end{equation*}
$$

satisfies the condition

$$
\begin{equation*}
\left|z^{2} f^{\prime}(z)+(1-\alpha) z f(z)+\alpha\right|<M_{n}(\alpha), \quad z \in U \tag{1.11}
\end{equation*}
$$

then

$$
\operatorname{Re}\left\{-\frac{z f^{\prime}(z)}{f(z)}\right\}>0
$$

## 2 Main results

Theorem 2.1 If $f \in \Sigma_{m}, m \in \mathbb{N}^{*}$, on the form

$$
f(z)=\frac{1}{z^{m}}+\sum_{k=m}^{\infty} a_{k} z^{k}
$$

and satisfies the condition

$$
\begin{equation*}
\left|(1-\alpha) m z^{m} f(z)+z^{m+1} f^{\prime}(z)+\alpha m\right|<M, \quad \alpha \in[0,2) \tag{2.1}
\end{equation*}
$$

then

$$
\begin{equation*}
\left|z^{m} f(z)-1\right|<\frac{M}{m(2-\alpha)} \tag{2.2}
\end{equation*}
$$

and this result is sharp.
Proof. If we let

$$
\begin{equation*}
p(z)=z^{m} f(z)-1 \tag{2.3}
\end{equation*}
$$

then $p \in \mathcal{H}_{2 m}$ and (2.1) can be rewritten as

$$
\begin{equation*}
\left|z p^{\prime}(z)-\alpha m p(z)\right|<M \tag{2.4}
\end{equation*}
$$

or

$$
\begin{equation*}
z p^{\prime}(z)-\alpha m p(z) \prec M z . \tag{2.5}
\end{equation*}
$$

If we take in Lemma 1.1

$$
q(z)=\frac{M z}{m(2-\alpha)}, \quad q \in \mathcal{H}(U)
$$

with $q(0)=0, q^{\prime}(0) \neq 0$ and

$$
\operatorname{Re}\left(1+\frac{z q^{\prime \prime}(z)}{q^{\prime}(z)}\right)>\frac{\alpha}{2}
$$

then from (1.7), $h(z)=M z$ and the result follows from Lemma 1.1, that is $p(z) \prec q(z)$

$$
z^{m} f(z)-1 \prec \frac{M z}{m(2-\alpha)}
$$

or

$$
\left|z^{m} f(z)-1\right|<\frac{M}{m(2-\alpha)} .
$$

Theorem 2.2 Let $m \in \mathbb{N}^{*}, 0 \leq \alpha<\frac{m}{m+1}$ and let

$$
\begin{equation*}
M(m, \alpha)=\frac{m(2-\alpha)(m-\alpha)}{m(1-\alpha)-\alpha+m(m-\alpha) \sqrt{\alpha^{2}+(2-\alpha)^{2}}} \tag{2.6}
\end{equation*}
$$

If $f \in \Sigma_{m}$ satisfies the condition

$$
\begin{equation*}
\left|(1-\alpha) m z^{m} f(z)+z^{m+1} f^{\prime}(z)+\alpha m\right|<M(m, \alpha) \tag{2.7}
\end{equation*}
$$

then $f \in \Omega_{m}(\alpha)$.
Proof. Let

$$
\begin{equation*}
0<M \leq M(m, \alpha) \tag{2.8}
\end{equation*}
$$

where $M(m, \alpha)$ is given by (2.6), and suppose that $f \in \Sigma_{m}$ satisfies the condition

$$
\begin{equation*}
\left|(1-\alpha) m z^{m} f(z)+z^{m+1} f^{\prime}(z)+\alpha m\right|<M . \tag{2.9}
\end{equation*}
$$

If we set

$$
\begin{equation*}
P(z)=z^{m} f(z) \tag{2.10}
\end{equation*}
$$

then by Theorem 2.1 we obtain

$$
\begin{equation*}
|P(z)-1|<\frac{M}{m(2-\alpha)} \equiv R, \quad z \in U \tag{2.11}
\end{equation*}
$$

From (2.6), we easily deduce $R<1$, which implies $P(z) \neq 0, z \in U$. Hence if we let

$$
\begin{equation*}
p(z)=-\alpha-\frac{z f^{\prime}(z)}{f(z)} \tag{2.12}
\end{equation*}
$$

then $p(z) \in \mathcal{H}[m-\alpha, 2 m]$ and (2.9) can be written in the form

$$
\begin{equation*}
|-p(z) P(z)+[m(1-\alpha)-\alpha] P(z)+\alpha m|<M \tag{2.13}
\end{equation*}
$$

We claim that this inequality implies $\operatorname{Re} p(z)>0, z \in U$. If this is false, then there exists a point $z_{0} \in U$, such that $p\left(z_{0}\right)=i \rho$, where $\rho$ is real. We will show that at such a point the negation of condition (2.13) holds, that is

$$
\begin{equation*}
\left|-i \rho P\left(z_{0}\right)+[m(1-\alpha)-\alpha] P\left(z_{0}\right)+\alpha m\right| \geq M \tag{2.14}
\end{equation*}
$$

for all real $\rho$.
If we let $P_{0}=P\left(z_{0}\right)$, one obtains

$$
\begin{aligned}
\mid-i \rho P_{0} & +[m(1-\alpha)-\alpha] P_{0}+\left.\alpha m\right|^{2}=\rho^{2}\left|P_{0}\right|^{2}+[m(1-\alpha)-\alpha]^{2}\left|P_{0}\right|^{2} \\
& +\alpha^{2} m^{2}+2 \alpha m[m(1-\alpha)-\alpha] \operatorname{Re} P_{0}+2 \alpha m \rho \operatorname{Im} P_{0}
\end{aligned}
$$

The inequality (2.14) is equivalent to

$$
\begin{equation*}
E \equiv \rho^{2}\left|P_{0}\right|^{2}+2 \alpha m \rho \operatorname{Im} P_{0}+[m(1-\alpha)-\alpha]^{2}\left|P_{0}\right|^{2}+\alpha^{2} m^{2}+ \tag{2.15}
\end{equation*}
$$

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$$
+2 \alpha m[m(1-\alpha)-\alpha] \operatorname{Re} P_{0}-R^{2} m^{2}(2-\alpha)^{2} \geq 0
$$

Since from (2.11) we have

$$
\left|P_{0}\right|>1-R \quad \text { and } \quad \operatorname{Re} P_{0}>1-R
$$

from (2.11) and (2.15) one obtains

$$
\begin{aligned}
& E \geq\left|P_{0}\right|^{2} \rho^{2}+2 \alpha m \operatorname{Im} P_{0} \rho+[m(1-\alpha)-\alpha]^{2}(1-R)^{2}+ \\
& +\alpha^{2} m^{2}+2 \alpha m[m(1-\alpha)-\alpha](1-R)-R^{2} m^{2}(2-\alpha)^{2} .
\end{aligned}
$$

Hence $E \geq 0$ if

$$
\begin{equation*}
\alpha^{2} m^{2}\left(\operatorname{Im} P_{0}\right)^{2} \leq\left|P_{0}\right|^{2}\left\{[(m(1-\alpha)-\alpha)(1-R)+\alpha m]^{2}-R^{2} m^{2}(2-\alpha)^{2}\right\} \tag{2.16}
\end{equation*}
$$

or

$$
\begin{equation*}
\alpha^{2} m^{2}\left(\operatorname{Im} P_{0}\right)^{2} \leq\left|P_{0}\right|^{2}\left\{[m-\alpha-[m(1-\alpha)-\alpha] R]^{2}-R^{2} m^{2}(2-\alpha)^{2}\right\} \tag{2.17}
\end{equation*}
$$

A simple geometric argument shows that the inequality (2.11) implies

$$
\begin{equation*}
\left(\operatorname{Im} P_{0}\right)^{2} \leq R^{2}\left|P_{0}\right|^{2} \tag{2.18}
\end{equation*}
$$

By comparing (2.17) and (2.18) we deduce that (2.14) holds if

$$
\begin{equation*}
\alpha^{2} m^{2} R^{2} \leq\{m-\alpha-[m(1-\alpha)-\alpha] R\}^{2}-R^{2} m^{2}(2-\alpha)^{2} \tag{2.19}
\end{equation*}
$$

or

$$
\begin{align*}
& R^{2}\left\{\alpha^{2} m^{2}+m^{2}(2-\alpha)^{2}-[m(1-\alpha)-\alpha]^{2}\right\}+  \tag{2.20}\\
& +2(m-\alpha)[m(1-\alpha)-\alpha] R-(m-\alpha)^{2} \leq 0
\end{align*}
$$

This last inequality holds if $R \leq R_{0}$, where

$$
\begin{equation*}
R_{0}=\frac{m-\alpha}{m(1-\alpha)-\alpha+m \sqrt{\alpha^{2}+(2-\alpha)^{2}}}, \quad 0 \leq \alpha<\frac{m}{m+1} \tag{2.21}
\end{equation*}
$$

that is $M \leq M(m, \alpha)$.
Thus we have a contradiction of (2.13), therefore $\operatorname{Re} p(z)>0, z \in U$ and $f \in$ $\Omega_{m}(\alpha)$.

Remark 2.1 Note that for the special case $m=1, \alpha=0$, the value $M(1,0)=2 / 3$ is the same with that obtained from (1.9) Lemma 1.2: $M_{1}(0)=2 / 3$.

We obtain the following criterion of starlikeness for meromorphic $m$-valent functions.

Corollary 2.1 Let $m \in \mathbb{N}^{*}$ and let $f \in \Sigma_{m}$ satisfies the condition

$$
\begin{equation*}
\left|m z^{m} f(z)+z^{m+1} f^{\prime}(z)\right|<\frac{2 m}{m+1} \tag{2.22}
\end{equation*}
$$

then $f \in \Omega_{m}^{*}$.

Since a function $f \in \Sigma_{m}$ can be written as

$$
\begin{equation*}
f(z)=\frac{1}{z^{m}}+g(z), \quad 0<|z|<1 \tag{2.23}
\end{equation*}
$$

where $g \in \mathcal{H}_{m}$, Theorem 2.2 can be rewritten in the following equivalent form, that is useful for the other results.
Corollary 2.2 Let $m \in \mathbb{N}^{*}, 0 \leq \alpha<\frac{m}{m+1}$ and let $f \in \Sigma_{m}$ have the form

$$
f(z)=\frac{1}{z^{m}}+g(z)
$$

where $g \in \mathcal{H}_{m}$. If

$$
\begin{equation*}
\left|(1-\alpha) m z^{m} g(z)+z^{m+1} g^{\prime}(z)\right|<M(m, \alpha), \quad z \in U \tag{2.24}
\end{equation*}
$$

where $M(m, \alpha)$ is given by (2.6), then $f \in \Omega_{m}(\alpha)$.
This form has an interesting interpretation in terms of integral operators. If we let

$$
\begin{equation*}
h(z)=(1-\alpha) m z^{m} g(z)+z^{m+1} g^{\prime}(z), \tag{2.25}
\end{equation*}
$$

then

$$
\begin{equation*}
g(z)=\frac{1}{z^{(1-\alpha) m}} \int_{0}^{z} h(t) t^{-(1+\alpha m)} d t \tag{2.26}
\end{equation*}
$$

which leads to the following result.
Corollary 2.3 Let $h \in \mathcal{H}_{2 m}$ and $M(m, \alpha)$ is given by (2.6) with $0 \leq \alpha<\frac{m}{m+1}$. If $h$ satisfies the condition

$$
\begin{equation*}
|h(z)| \leq M(m, \alpha), \quad z \in U \tag{2.27}
\end{equation*}
$$

then

$$
\begin{equation*}
f(z)=\frac{1}{z}+\frac{1}{z^{(1-\alpha) m}} \int_{0}^{z} h(t) t^{-(1+\alpha m)} d t \in \Omega_{m}(\alpha) \tag{2.28}
\end{equation*}
$$

Example 2.1 For the Corollary 2.3 we consider the following function

$$
\begin{equation*}
h(z)=a z^{3}(z-\sin z) \tag{2.29}
\end{equation*}
$$

Since $h \in \mathcal{H}_{6}$ we deduce that $m=3$ and we choose for $\alpha$ a value such that $0 \leq \alpha<\frac{m}{m+1}$. Let the value be $\alpha=\frac{2}{3}$. Then, from (2.28) we get

$$
\begin{equation*}
f(z)=\frac{1}{z}+\frac{a}{z} \int_{0}^{z} t^{3}(t-\sin t) t^{-3} d t \tag{2.30}
\end{equation*}
$$

or

$$
\begin{equation*}
f(z)=\frac{1}{z}\left(1-2 a \sin ^{2} \frac{z}{2}\right)+\frac{a z}{2} \tag{2.31}
\end{equation*}
$$

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From (2.27) we obtain

$$
\begin{equation*}
|h(z)| \leq M(m, \alpha)=M\left(3, \frac{2}{3}\right) \tag{2.32}
\end{equation*}
$$

The above inequality leads to the relation

$$
\begin{equation*}
\left|a z^{3}\right||z-\sin z| \leq|a| \frac{e^{2}+2 e-1}{2 e} \tag{2.33}
\end{equation*}
$$

The condition (2.32) will be satisfied if

$$
\begin{equation*}
|a| \frac{e^{2}+2 e-1}{2 e} \leq \frac{21}{1+7 \sqrt{20}} \tag{2.34}
\end{equation*}
$$

and we obtain

$$
|a| \leq \frac{42 e}{\left(2 e+e^{2}-1\right)(1+7 \sqrt{20})}=0.298 \ldots
$$

Hence, if we take $a=\frac{1}{4}$ we conclude that

$$
f(z)=\frac{1}{z}\left(1-\frac{1}{2} \sin ^{2} \frac{z}{2}\right)+\frac{z}{8} \in \Omega_{3}\left(\frac{2}{3}\right) .
$$

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# Generalized classes of uniformly convex functions 

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Submitted by: Jan Stankiewicz


#### Abstract

In this paper we introduce some subclasses of analytic functions with varying argument of coefficients. These classes are defined in terms of the Hadamard product and generalize the well-known classes of uniformly convex functions. We investigate the coefficients estimates, distortion properties, radii of starlikeness and convexity for defined classes of functions


## AMS Subject Classification: 30C45

Key Words and Phrases: analytic functions, varying arguments, subordination, Hadamard product

## 1 Introduction

Let $\mathcal{A}$ denote the class of functions which are analytic in $\mathcal{U}=\mathcal{U}(1)$, where

$$
\mathcal{U}(r)=\{z \in \mathbb{C}:|z|<r\} .
$$

and let $\mathcal{A}(p, k)(p, k \in \mathbb{N}=\{1,2,3 \ldots\}, p<k)$ denote the class of functions $f \in \mathcal{A}$ of the form

$$
\begin{equation*}
f(z)=a_{p} z^{p}+\sum_{n=k}^{\infty} a_{n} z^{n} \quad\left(z \in \mathcal{U} ; a_{p}>0\right) . \tag{1}
\end{equation*}
$$

For multivalent fuction $f \in \mathcal{A}(p, k)$ the normalization

$$
\begin{equation*}
\left.\frac{f(z)}{z^{p-1}}\right|_{z=0}=0 \text { and }\left.\frac{f(z)}{z^{p}}\right|_{z=0}=1 . \tag{2}
\end{equation*}
$$

is clasical. One can obtain interesting results by applying normalization of the form

$$
\begin{equation*}
\left.\frac{f(z)}{z^{p-1}}\right|_{z=0}=0 \text { and }\left.\frac{f(z)}{z^{p}}\right|_{z=\rho}=1 \tag{3}
\end{equation*}
$$

where $\rho$ is a fixed point of the unit $\operatorname{disk} \mathcal{U}$. In particular, for $p=1$ we obtain Montel's normaliztion ( $c f$. [1]). We see that for $\rho=0$ the normalization (3) is the clasical.

We denote by $\mathcal{A}_{\rho}(p, k)$ the classes of functions $f \in \mathcal{A}(p, k)$ with the normalization (3). It will be called the class of functions with two fixed points.

Also, by $\mathcal{T}(p, k ; \eta) \quad(\eta \in \mathbb{R})$ we denote the class of functions $f \in \mathcal{A}(p, k)$ of the form (1) for which

$$
\begin{equation*}
\arg \left(a_{n}\right)=\pi+(p-n) \eta \quad(n=k, k+1, \ldots) \tag{4}
\end{equation*}
$$

For $\eta=0$ we obtain the class $\mathcal{T}(p, k ; 0)$ of functions with negative coefficients. Moreover, we define

$$
\begin{equation*}
\mathcal{T}(p, k):=\bigcup_{\eta \in \mathbb{R}} \mathcal{T}(p, k ; \eta) \tag{5}
\end{equation*}
$$

The classes $\mathcal{T}(p, k)$ and $\mathcal{T}(p, k ; \eta)$ are called the classes of functions with varying argument of coefficients. The class $\mathcal{T}(1,2)$ was introduced by Silverman [2] (see also [3]).

Let $\alpha \in\langle 0, p), r \in(0,1\rangle$. A function $f \in \mathcal{A}(p, k)$ is said to be convex of order $\alpha$ in $\mathcal{U}(r)$ if and only if

$$
\operatorname{Re}\left(1+\frac{z f^{\prime \prime}(z)}{f^{\prime}(z)}\right)>\alpha(z \in \mathcal{U}(r))
$$

A function $f \in \mathcal{A}(p, k)$ is said to be starlike of order $\alpha$ in $\mathcal{U}(r)$ if and only if

$$
\begin{equation*}
\operatorname{Re}\left(\frac{z f^{\prime}(z)}{f(z)}\right)>\alpha \quad(z \in \mathcal{U}(r)) \tag{6}
\end{equation*}
$$

We denote by $\mathcal{S}^{c}(\alpha)$ the class of all functions $f \in \mathcal{A}(p, p+1)$, which are convex of order $\alpha$ in $\mathcal{U}$ and by $\mathcal{S}_{p}^{*}(\alpha)$ we denote the class of all functions $f \in \mathcal{A}(p, p+1)$, which are starlike of order $\alpha$ in $\mathcal{U}$. We also set

$$
\mathcal{S}^{c}=\mathcal{S}_{1}^{c}(0) \text { and } \mathcal{S}^{*}=\mathcal{S}_{1}^{*}(0)
$$

It is easy to show that for a function $f$ from the class $\mathcal{T}(p, k)$ the condition (6) is equivalent to the following

$$
\begin{equation*}
\left|\frac{z f^{\prime}(z)}{f(z)}-p\right|<p-\alpha \quad(z \in \mathcal{U}(r)) \tag{7}
\end{equation*}
$$

Let $\mathcal{B}$ be a subclass of the class $\mathcal{A}(p, k)$. We define the radius of starlikeness of order $\alpha$ and the radius of convexity of order $\alpha$ for the class $\mathcal{B}$ by

$$
\begin{aligned}
R_{\alpha}^{*}(\mathcal{B}) & =\inf _{f \in \mathcal{B}}(\sup \{r \in(0,1]: f \text { is starlike of order } \alpha \text { in } \mathcal{U}(r)\}), \\
R_{\alpha}^{c}(\mathcal{B}) & =\inf _{f \in \mathcal{B}}(\sup \{r \in(0,1]: f \text { is convex of order } \alpha \text { in } \mathcal{U}(r)\}),
\end{aligned}
$$

respectively.

We say that a function $f \in \mathcal{A}$ is subordinate to a function $F \in \mathcal{A}$, and write $f(z) \prec$ $F(z)$ (or simply $f \prec F$ ), if and only if there exists a function $\omega \in \mathcal{A}(|\omega(z)| \leq|z|, \quad z \in \mathcal{U})$ such that

$$
f(z)=F(\omega(z)) \quad(z \in \mathcal{U})
$$

In particular, if $F$ is univalent in $\mathcal{U}$, we have the following equivalence.

$$
f(z) \prec F(z) \Longleftrightarrow f(0)=F(0) \text { and } f(\mathcal{U}) \subset F(\mathcal{U})
$$

For functions $f, g \in \mathcal{A}$ of the form

$$
f(z)=\sum_{n=0}^{\infty} a_{n} z^{n} \text { and } g(z)=\sum_{n=0}^{\infty} b_{n} z^{n}
$$

by $f * g$ we denote the Hadamard product (or convolution) of $f$ and $g$, defined by

$$
(f * g)(z)=\sum_{n=0}^{\infty} a_{n} b_{n} z^{n} \quad(z \in \mathcal{U})
$$

Let $\gamma, \delta$ be real parameters, $0 \leq \gamma<1, \delta \geq 0$, and let $\varphi, \phi \in \mathcal{A}_{0}(p, k)$.
By $\mathcal{W}(p, k ; \phi, \varphi ; \gamma, \delta)$ we denote the class of functions $f \in \mathcal{A}(p, k)$ such that

$$
\begin{equation*}
(\varphi * f)(z) \neq 0 \quad(z \in \mathcal{U} \backslash\{0\}) \tag{8}
\end{equation*}
$$

and

$$
\begin{equation*}
\operatorname{Re}\left\{\frac{(\phi * f)(z)}{(\varphi * f)(z)}-\gamma\right\}>\delta\left|\frac{(\phi * f)(z)}{(\varphi * f)(z)}-1\right| \quad(z \in \mathcal{U}) \tag{9}
\end{equation*}
$$

Also, let us denote

$$
\begin{array}{rll}
\mathcal{T} \mathcal{W}(p, k ; \phi, \varphi ; \gamma, \delta) & : & =\mathcal{T}(p, k) \cap \mathcal{W}(p, k ; \phi, \varphi ; \gamma, \delta), \\
\mathcal{T} \mathcal{W}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta) & : & =\mathcal{T}(p, k ; \eta) \cap \mathcal{W}(p, k ; \phi, \varphi ; \gamma, \delta), \\
\mathcal{T} \mathcal{W}_{\rho}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta) & :=\mathcal{A}_{\rho}(p, k) \cap \mathcal{T} \mathcal{W}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta), \\
\mathcal{T} \mathcal{W}_{\rho}(p, k ; \phi, \varphi ; \gamma, \delta) & : & =\mathcal{A}_{\rho}(p, k) \cap \mathcal{T} \mathcal{W}(p, k ; \phi, \varphi ; \gamma, \delta)
\end{array}
$$

For the presented investigations we assume that $\varphi, \phi$ are the functions of the form

$$
\begin{equation*}
\varphi(z)=z^{p}+\sum_{n=k}^{\infty} \alpha_{n} z^{n}, \quad \phi(z)=z^{p}+\sum_{n=k}^{\infty} \beta_{n} z^{n} \quad(z \in \mathcal{U}) \tag{10}
\end{equation*}
$$

where

$$
0 \leq \alpha_{n}<\beta_{n} \quad(n=k, k+1, \ldots)
$$

Moreover, let us put

$$
\begin{equation*}
d_{n}:=(\delta+1) \beta_{n}-(\delta+\gamma) \alpha_{n} \quad(n=k, k+1, \ldots) \tag{11}
\end{equation*}
$$

The families $\mathcal{W}_{\rho}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)$ and $\mathcal{W}_{\rho}(p, k ; \phi, \varphi ; \gamma, \delta)$ unify various new and well-known classes of analytic functions. In particular, the class

$$
\mathcal{W}_{\rho}(\varphi ; \gamma, \delta ; \eta):=\mathcal{W}_{\rho}\left(p, k ; \frac{z \varphi^{\prime}(z)}{p}, \varphi(z) ; \gamma, \delta ; \eta\right)
$$

contains functions $f \in \mathcal{A}_{\rho}(p, k)$, such that

$$
\operatorname{Re}\left\{\frac{z(\varphi * f)^{\prime}(z)}{p(\varphi * f)(z)}-\gamma\right\}>\delta\left|\frac{z(\varphi * f)^{\prime}(z)}{p(\varphi * f)(z)}-1\right| \quad(z \in \mathcal{U}) .
$$

The class

$$
\mathcal{H}_{\mathcal{T}}(\varphi ; \gamma, \delta):=\mathcal{T} \mathcal{W}_{0}(1,2 ; \varphi ; \gamma, \delta ; 0)
$$

was introduced and studied by Raina and Bansal [4]. If we set

$$
h\left(\alpha_{1}, z\right):=z_{q} F_{s}\left(\alpha_{1}, \ldots, \alpha_{q} ; \beta_{1}, \ldots, \beta_{s} ; z\right),
$$

where ${ }_{q} F_{s}$ is the generalized hypergeometric function (see for details [5] and [6]), then we obtain the class

$$
\mathcal{U H}(q, s, \lambda, \gamma, \delta):=\mathcal{T} \mathcal{W}_{0}\left(1,2 ; \lambda h\left(\alpha_{1}+1, z\right)+(1-\lambda) h\left(\alpha_{1}, z\right) ; \gamma, \delta ; 0\right) \quad(0 \leq \lambda \leq 1)
$$

defined by Srivastava et al. [7]. The classes

$$
\begin{aligned}
\delta-\operatorname{UST}(\gamma) & =\mathcal{W}_{0}\left(1,2 ; \frac{z}{1-z} ; \gamma, \delta\right) \\
\delta-U C V(\gamma) & =\mathcal{W}_{0}\left(1,2 ; \frac{z}{(1-z)^{2}} ; \gamma, \delta\right)
\end{aligned}
$$

are the well-known classes of of $\delta$-starlike function of order $\gamma$ and $\delta$-uniformly convex function of order $\gamma$, respectively. In particular, the classes $U C V:=\operatorname{UCV}(1,0)$, $\delta-U C V:=U C V(\delta, 0)$ were introduced by Goodman [8] (see also [9, 10, 11]), and Kanas and Wisniowska [12], respectively.

Many other classes, are also particular cases of the class investigated here, see for example [13, 14, 15].

The object of the present paper is to investigate the coefficients estimates, distortion properties and the radii of starlikeness and convexity.

## 2 Coefficients estimates

We first mention a sufficient condition for the function to belong to the class $\mathcal{W}(p, k ; \phi, \varphi ; \gamma, \delta)$.
Theorem 1 Let $\left\{d_{n}\right\}$ be defined by (11), and let $0 \leq \gamma<1$. If a function $f$ of the form (1) satisfies the condition

$$
\begin{equation*}
\sum_{n=k}^{\infty} d_{n}\left|a_{n}\right| \leq(1-\gamma) a_{p} \tag{12}
\end{equation*}
$$

then $f$ belongs to the class $\mathcal{W}(p, k ; \phi, \varphi ; \gamma, \delta)$.

Proof. By definition of the class $\mathcal{W}(p, k ; \phi, \varphi ; \gamma, \delta)$, it suffices to show that

$$
\begin{equation*}
\delta\left|\frac{(\phi * f)(z)}{(\varphi * f)(z)}-1\right|-\operatorname{Re}\left\{\frac{(\phi * f)(z)}{(\varphi * f)(z)}-1\right\} \leq 1-\gamma \quad(z \in \mathcal{U}) \tag{13}
\end{equation*}
$$

Simply calculations give

$$
\begin{aligned}
& \delta\left|\frac{(\phi * f)(z)}{(\varphi * f)(z)}-1\right|-\operatorname{Re}\left\{\frac{(\phi * f)(z)}{(\varphi * f)(z)}-1\right\} \\
\leq & (\delta+1)\left|\frac{(\phi * f)(z)}{(\varphi * f)(z)}-1\right| \leq(\delta+1) \frac{\sum_{n=k}^{\infty}\left(\beta_{n}-\alpha_{n}\right)\left|a_{n}\right||z|^{n-p}}{a_{p}-\sum_{n=k}^{\infty} \alpha_{n}\left|a_{n}\right||z|^{n-p}}
\end{aligned}
$$

Now the last expression is bounded above by $(1-\gamma)$ if (12) holds. Whence $f \in$ $\mathcal{W}(p, k ; \phi, \varphi ; \gamma, \delta)$.

Our next theorem shows that the condition (12) is necessary as well for functions of the form (1), with (4) to belong to the class $\mathcal{T} \mathcal{W}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)$.

Theorem 2 Let $f$ be a function of the form (1), satisfying the argument property (4). Then $f$ belongs to the class $\mathcal{T} \mathcal{W}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)$ if and only if the condition (12) holds true.

Proof. In view of Theorem 1 we need only show that each function $f$ from the class $\mathcal{T} \mathcal{W}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)$ satisfies the coefficient inequality (12). Let a function $f$ of the form (1), satisfying the argument property (4) belong to the class $\mathcal{T} \mathcal{W}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)$. Then by (9), we have

$$
\delta\left|\frac{a_{p} z^{p}+\sum_{n=k}^{\infty} \beta_{n} a_{n} z^{n}}{a_{p} z^{p}+\sum_{n=k}^{\infty} \alpha_{n} a_{n} z^{n}}-1\right|<\operatorname{Re}\left\{\frac{a_{p} z^{p}+\sum_{n=k}^{\infty} \beta_{n} a_{n} z^{n}}{a_{p} z^{p}+\sum_{n=k}^{\infty} \alpha_{n} a_{n} z^{n}}-\gamma\right\}
$$

or equivalently

$$
\delta\left|\frac{\sum_{n=k}^{\infty}\left(\beta_{n}-\alpha_{n}\right) a_{n} z^{n-p}}{a_{p}+\sum_{n=k}^{\infty} \alpha_{n} a_{n} z^{n-p}}\right|<\operatorname{Re}\left\{\frac{(1-\gamma) a_{p}+\sum_{n=k}^{\infty}\left(\beta_{n}-\gamma \alpha_{n}\right) a_{n} z^{n-p}}{a_{p}+\sum_{n=k}^{\infty} \alpha_{n} a_{n} z^{n-p}}\right\} .
$$

In view of (4), we set $z=r e^{i \eta}(0 \leq r<1)$ in the above inequality to obtain

$$
\frac{\sum_{n=k}^{\infty} \delta\left(\beta_{n}-\alpha_{n}\right)\left|a_{n}\right| r^{n-p}}{a_{p}-\sum_{n=k}^{\infty} \alpha_{n}\left|a_{n}\right| r^{n-p}}<\frac{(1-\gamma) a_{p}-\sum_{n=k}^{\infty}\left(\beta_{n}-\gamma \alpha_{n}\right)\left|a_{n}\right| r^{n-p}}{a_{p}-\sum_{n=k}^{\infty} \alpha_{n}\left|a_{n}\right| r^{n-p}}
$$

Thus, by (8) we have

$$
\sum_{n=k}^{\infty}\left[(\delta+1) \beta_{n}-(\delta+\gamma) \alpha_{n}\right]\left|a_{n}\right| r^{n-p}<(1-\gamma) a_{p}
$$

which, upon letting $r \rightarrow 1^{-}$, readily yields the assertion (12).
By applying Theorem 2 we can deduce following result.
Theorem 3 Let $f$ be a function of the form (1), satisfying the argument property (4). A function $f$ of the form (1) belongs to the class $\mathcal{T} \mathcal{W}_{\rho}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)$ if and only if it satisfies (3) and

$$
\begin{equation*}
\sum_{n=k}^{\infty}\left(d_{n}-(1-\gamma)|\rho|^{n-p}\right)\left|a_{n}\right| \leq 1-\gamma \tag{14}
\end{equation*}
$$

where $\left\{d_{n}\right\}$ is defined by (11).

Proof. For a function $f$ of the form (1) with the normalization (3), we have

$$
\begin{equation*}
a_{p}=1+\sum_{n=k}^{\infty}\left|a_{n}\right||\rho|^{n-p} \tag{15}
\end{equation*}
$$

Applying the equality (15) to (12), we obtain the assertions (14).
Since the condition (14) is independent of $\eta$, Theorem 3 yields the following theorem.

Theorem 4 Let $f$ be a function of the form (1), satisfying the argument property (4). Then $f$ belongs to the class $\mathcal{T} \mathcal{W}_{\rho}(p, k ; \phi, \varphi ; \gamma, \delta)$ if and only if the condition (14) holds true.

By applying Theorem 3 we obtain the following lemma.
Lemma 1 Let $\left\{d_{n}\right\}$ be defined by (11), $\rho \in \mathcal{U}$, and let us assume, that there exists an integer $n_{0}\left(n_{0} \in \mathbb{N}_{k}=\{k, k+1, \ldots\}\right)$ such that

$$
\begin{equation*}
d_{n_{0}}-(1-\gamma)|\rho|^{n_{0}-p} \leqq 0 \tag{16}
\end{equation*}
$$

Then the function

$$
f_{n_{0}}(z)=\left(1+a \rho^{n_{0}-p}\right) z^{p}-a e^{i\left(p-n_{0}\right) \eta} z^{n_{0}}
$$

belongs to the class $\mathcal{T} \mathcal{W}_{\rho}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)$ for all positive real numbers a. Moreover, for all $n\left(n \in \mathbb{N}_{k}\right)$ such that

$$
\begin{equation*}
d_{n}-(1-\gamma)|\rho|^{n-p}>0 \tag{17}
\end{equation*}
$$

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the functions

$$
f_{n}(z)=\left(1+a \rho^{n_{0}-p}+b z^{n-p}\right) z^{p}-a e^{i\left(p-n_{0}\right) \eta} z^{n_{0}}-b e^{i(p-n) \eta} z^{n}
$$

where

$$
b=\frac{1-\gamma+\left((1-\gamma)|\rho|^{n_{0}-p}-d_{n_{0}}\right) a}{d_{n}-(1-\gamma)|\rho|^{n-p}}
$$

belong to the class $\mathcal{T} \mathcal{W}_{\rho}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)$.
By Lemma 1 and Theorem 3, we have following corollary.
Corollary 1 Let a function $f$ of the form (1) belongs to the class
$\mathcal{T} \mathcal{W}_{\rho}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)$ and let $\left\{d_{n}\right\}$ be defined by (11). Then all of the coefficients $a_{n}$ for which

$$
d_{n}-(1-\gamma)|\rho|^{n-p}=0
$$

are unbounded. Moreover, if there exists an integer $n_{0}\left(n_{0} \in \mathbb{N}_{k}=\{k, k+1, \ldots\}\right)$ such that

$$
d_{n_{0}}-(1-\gamma)|\rho|^{n_{0}-p}<0
$$

then all of the coefficients of the function $f$ are unbounded. In the remaining cases

$$
\left|a_{n}\right| \leqq \frac{1-\gamma}{d_{n}-(1-\gamma)|\rho|^{n-p}}
$$

The result is sharp, the functions $f_{n}$ of the form

$$
f_{n, \eta}(z)=\frac{d_{n} z^{p}-(1-\gamma) e^{i(p-n) \eta} z^{n}}{d_{n}-(1-\gamma)|\rho|^{n-p}} \quad(z \in \mathcal{U} ; n=k, k+1, \ldots)
$$

are the extremal functions.
Remark 1 The coefficients estimates for the class $\mathcal{T} \mathcal{W}_{\rho}(p, k ; \phi, \varphi ; \gamma, \delta)$ are the same as for the class $\mathcal{T} \mathcal{W}_{\rho}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)$.

By puting $\rho=0$ in Theorems 3 and 4, and Corollary 1, we have the corollaries listed below.

Corollary 2 Let $f$ be a function of the form (1), satisfying the argument property (4). A function $f$ of the form (1) belongs to the class $\mathcal{T} \mathcal{W}_{0}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)$ if and only if

$$
\begin{equation*}
\sum_{n=k}^{\infty} d_{n}\left|a_{n}\right| \leq 1-\gamma \tag{18}
\end{equation*}
$$

where $\left\{d_{n}\right\}$ is defined by (11).

Corollary 3 Let $f$ be a function of the form (1), satisfying the argument property (4). Then $f$ belongs to the class $\mathcal{T} \mathcal{W}_{0}(p, k ; \phi, \varphi ; \gamma, \delta)$ if and only if the condition (18) holds true.

Corollary 4 If a function $f$ of the form (1) belongs to the class $\mathcal{T} \mathcal{W}_{0}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)$, then

$$
\begin{equation*}
\left|a_{n}\right| \leq \frac{1-\gamma}{d_{n}} \quad(n=k, k+1, \ldots) \tag{19}
\end{equation*}
$$

where $d_{n}$ is defined by (11). The result is sharp. The functions $f_{n, \eta}$ of the form

$$
\begin{equation*}
f_{n, \eta}(z)=z^{p}-\frac{1-\gamma}{d_{n}} e^{i(p-n) \eta} z^{n} \quad(z \in \mathcal{U} ; n=k, k+1, \ldots) \tag{20}
\end{equation*}
$$

are the extremal functions.
Corollary 5 If a function $f$ of the form (1) belongs to the class $\mathcal{T} \mathcal{W}_{0}(p, k ; \phi, \varphi ; \gamma, \delta)$, then the coefficients estimates (19) holds true. The result is sharp. The functions $f_{n, \eta}$ $(\eta \in \mathbb{R})$ of the form (20) are the extremal functions.

## 3 Distortion theorems

From Theorem 2 we have the following lemma.
Lemma 2 Let a function $f$ of the form (1) belong to the class $\mathcal{T W}_{\rho}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)$. If the sequence $\left\{d_{n}\right\}$ defined by (11) satisfies the inequality

$$
\begin{equation*}
0<d_{k}-(1-\gamma)|\rho|^{k-p} \leq d_{n}-(1-\gamma)|\rho|^{n-p} \quad(n=k, k+1, \ldots) \tag{21}
\end{equation*}
$$

then

$$
\sum_{n=k}^{\infty}\left|a_{n}\right| \leq \frac{1-\gamma}{d_{k}-(1-\gamma)|\rho|^{k-p}}
$$

Moreover, if

$$
\begin{equation*}
0<\frac{d_{k}-(1-\gamma)|\rho|^{k-p}}{k} \leq \frac{d_{n}-(1-\gamma)|\rho|^{n-p}}{n} \quad(n=k, k+1, \ldots) \tag{22}
\end{equation*}
$$

then

$$
\sum_{n=k}^{\infty} n\left|a_{n}\right| \leq \frac{k(1-\gamma)}{d_{k}-(1-\gamma)|\rho|^{k-p}}
$$

Remark 2 The second part of Lemma 2 we can rewritten in terms of $\sigma$-neighborhood $N_{\sigma}$ defined by

$$
N_{\sigma}=\left\{f(z)=a_{p} z^{p}+\sum_{n=k}^{\infty} a_{n} z^{n} \in \mathcal{T}(p, k ; \eta): \quad \sum_{n=k}^{\infty} n\left|a_{n}\right| \leq \sigma\right\}
$$

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in the following form:
if the sequence $\left\{d_{n}\right\}$ defined by (11) satisfies (22), then

$$
\mathcal{T} \mathcal{W}_{\rho}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta) \subset N_{\sigma}
$$

where

$$
\delta=\frac{k(1-\gamma)}{d_{k}-(1-\gamma)|\rho|^{k-p}}
$$

Theorem 5 Let a function $f$ belong to the class $\mathcal{T} \mathcal{W}_{\rho}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)$ and let $|z|=$ $r<1$. If the sequence $\left\{d_{n}\right\}$ defined by (11) satisfies (21), then

$$
\begin{equation*}
\phi(r) \leq|f(z)| \leq \frac{d_{k} r^{p}+(1-\gamma) r^{k}}{d_{k}-(1-\gamma)|\rho|^{k-p}} \tag{23}
\end{equation*}
$$

where

$$
\phi(r):=\left\{\begin{array}{cc}
r^{p} & (r \leqq \rho)  \tag{24}\\
\frac{d_{k} r^{p}-(1-\gamma) r^{k}}{d_{k}-(1-\gamma)|\rho|^{k-p}} & (r>\rho) .
\end{array}\right.
$$

Moreover, if (22) holds, then

$$
\begin{equation*}
p a_{p} r^{p-1}-\frac{k(1-\gamma)}{d_{k}-(1-\gamma)|\rho|^{k-p}} r^{k-1} \leq\left|f^{\prime}(z)\right| \leq \frac{p d_{k} r^{p}+k(1-\gamma) r^{k-1}}{d_{k}-(1-\gamma)|\rho|^{k-p}} \tag{25}
\end{equation*}
$$

The result is sharp, with the extremal function $f_{k, \eta}$ of the form (20) and $f(z)=z$.
Proof. Suppose that the function $f$ of the form (1) belongs to the class $\mathcal{T} \mathcal{W}_{\rho}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)$. By Lemma 2 we have

$$
\begin{aligned}
|f(z)| & =\left|a_{p} z^{p}+\sum_{n=k}^{\infty} a_{n} z^{n}\right| \leq r^{p}\left(a_{p}+\sum_{n=k}^{\infty}\left|a_{n}\right| r^{n-p}\right) \\
& \leq r^{p}\left(1+\sum_{n=k}^{\infty}\left|a_{n}\right||\rho|^{n-p}+\sum_{n=k}^{\infty}\left|a_{n}\right| r^{n-p}\right) \\
& \leq r^{p}\left(1+\left(|\rho|^{k-p}+r^{k-p}\right) \sum_{n=k}^{\infty}\left|a_{n}\right|\right) \leq \frac{d_{k} r^{p}+(1-\gamma) r^{k}}{d_{k}-(1-\gamma)|\rho|^{k-p}},
\end{aligned}
$$

and

$$
\begin{equation*}
|f(z)| \geq r^{p}\left(a_{p}-\sum_{n=k}^{\infty}\left|a_{n}\right| r^{n-p}\right)=r^{p}\left(1+\sum_{n=k}^{\infty}\left(|\rho|^{n-p}-r^{n-p}\right)\left|a_{n}\right|\right) \tag{26}
\end{equation*}
$$

If $r \leqq \rho$, then we obtain $|f(z)| \geqq r^{p}$. If $r>\rho$, then the sequence $\left\{\left(\rho^{n-1}-r^{n-1}\right)\right\}$ is decreasing and negative. Thus, by (26), we obtain

$$
|f(z)| \geq r^{p}\left(1-\left(r^{k-p}-|\rho|^{k-p}\right) \sum_{n=2}^{\infty} a_{n}\right) \geqq \frac{d_{k} r^{p}-(1-\gamma) r^{k}}{d_{k}-(1-\gamma)|\rho|^{k-p}}
$$

and we have the assertion (23). Making use of Lemma 2, in conjunction with (15), we readily obtain the assertion (25) of Theorem 5.

Theorem 5 implies the following results.

Corollary 6 Let a function $f$ belong to the class $\mathcal{T} \mathcal{W}_{\rho}(p, k ; \phi, \varphi ; \gamma, \delta)$. If the sequence $\left\{d_{n}\right\}$ defined by (11) satisfies (21), then the assertion (23) holds true. Moreover, if we assume (22), then then the assertion (25) holds true. The result is sharp, with the extremal functions $f_{k, \eta}(\eta \in \mathbb{R})$ of the form (20).
Corollary 7 Let a function $f$ belong to the class $\mathcal{T} \mathcal{W}_{0}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)$ and let the sequence $\left\{d_{n}\right\}$ be defined by (11). If

$$
\begin{equation*}
d_{k} \leq d_{n} \quad(n=k, k+1, \ldots) \tag{27}
\end{equation*}
$$

then

$$
\begin{equation*}
r^{p}-\frac{1-\gamma}{d_{k}} r^{k} \leq|f(z)| \leq r^{p}+\frac{1-\gamma}{d_{k}} r^{k} \quad(|z|=r<1) . \tag{28}
\end{equation*}
$$

Moreover, if

$$
\begin{equation*}
n d_{k} \leq k d_{n} \quad(n=k, k+1, \ldots) \tag{29}
\end{equation*}
$$

then

$$
\begin{equation*}
p r^{p-1}-\frac{k(1-\gamma)}{d_{k}} r^{k-1} \leq\left|f^{\prime}(z)\right| \leq p r^{p-1}+\frac{k(1-\gamma)}{d_{k}} r^{k-1} \quad(|z|=r<1) \tag{30}
\end{equation*}
$$

The result is sharp, with the extremal function $f_{k, \eta}$ of the form (20).
Corollary 8 Let a function $f$ belong to the class $\mathcal{T} \mathcal{W}_{0}(p, k ; \phi, \varphi ; \gamma, \delta)$. If the sequence $\left\{d_{n}\right\}$ defined by (11) satisfies (27), then the assertion (28) holds true. Moreover, if we assume (29), then then the assertion (28) holds true. The result is sharp, with the extremal functions $f_{k, \eta} \quad(\eta \in \mathbb{R})$ of the form (20).

## 4 The Radii of convexity and starlikeness

Theorem 6 The radius of starlikeness of order $\alpha$ for the class $\mathcal{T} \mathcal{W}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)$ is given by

$$
\begin{equation*}
R_{\alpha}^{*}(\mathcal{T W}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta))=\inf _{n \geq k}\left(\frac{(p-\alpha) d_{n}}{(n-\alpha)(1-\gamma)}\right)^{\frac{1}{n-p}} \tag{31}
\end{equation*}
$$

where $d_{n}$ is defined by (11). The functions $f_{n, \eta}$ of the form

$$
\begin{equation*}
f_{n, \eta}(z)=a_{p}\left(z^{p}-\frac{1-\gamma}{d_{n}} e^{i(p-n) \eta} z^{n}\right) \quad\left(z \in \mathcal{U} ; n=k, k+1, \ldots ; a_{p}>0\right) \tag{32}
\end{equation*}
$$

are the extremal functions.
Proof. A function $f \in \mathcal{T}(p, k ; \eta)$ of the form (1) is starlike of order $\alpha$ in the disk $\mathcal{U}(r), 0<r \leq 1$, if and only if it satisfies the condition (7). Since

$$
\left|\frac{z f^{\prime}(z)}{f(z)}-p\right|=\left|\frac{\sum_{n=k}^{\infty}(n-p) a_{n} z^{n}}{a_{p} z^{p}+\sum_{n=k}^{\infty} a_{n} z^{n}}\right| \leq \frac{\sum_{n=k}^{\infty}(n-p)\left|a_{n}\right||z|^{n-p}}{a_{p}-\sum_{n=k}^{\infty}\left|a_{n}\right||z|^{n-p}}
$$

putting $|z|=r$ the condition (7) is true if

$$
\begin{equation*}
\sum_{n=k}^{\infty} \frac{n-\alpha}{p-\alpha}\left|a_{n}\right| r^{n-p} \leq a_{p} \tag{33}
\end{equation*}
$$

By Theorem 2, we have

$$
\begin{equation*}
\sum_{n=k}^{\infty} \frac{d_{n}}{1-\gamma}\left|a_{n}\right| \leq a_{p} \tag{34}
\end{equation*}
$$

Thus, the condition (33) is true if

$$
\frac{n-\alpha}{p-\alpha} r^{n-p} \leq \frac{d_{n}}{1-\gamma} \quad(n=k, k+1, \ldots)
$$

that is, if

$$
\begin{equation*}
r \leq\left(\frac{(p-\alpha) d_{n}}{(n-\alpha)(1-\gamma)}\right)^{\frac{1}{n-p}} \quad(n=k, k+1, \ldots) \tag{35}
\end{equation*}
$$

It follows that each function $f \in \mathcal{T} \mathcal{W}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)$ is starlike of order $\alpha$ in the disk $\mathcal{U}(r)$, where

$$
r=\inf _{n \geq k}\left(\frac{(p-\alpha) d_{n}}{(n-\alpha)(1-\gamma)}\right)^{\frac{1}{n-p}}
$$

The functions $f_{n, \eta}$ of the form (32) realize equality in (34), and the radius $r$ can not be larger. Thus we have (31).

The following result may be proved in much the same way as Theorem 6.
Theorem 7 The radius of convexity of order $\alpha$ for the class $\mathcal{T} \mathcal{W}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)$ is given by

$$
R_{\alpha}^{c}(\mathcal{T} \mathcal{W}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta))=\inf _{n \geq k}\left(\frac{(p-\alpha) d_{n}}{n(n-\alpha)(1-\gamma)}\right)^{\frac{1}{n-p}}
$$

where $d_{n}$ is defined by (11). The functions $f_{n, \eta}$ of the form (32) are the extremal functions.

It is clear that for

$$
a_{p}=\frac{d_{n}}{d_{n}-(1-\gamma)|\rho|^{n-p}}>0
$$

the extremal functions $f_{n, \eta}$ of the form (32) belong to the class $\mathcal{T} \mathcal{W}_{\rho}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)$. Moreover, we have

$$
\mathcal{T} \mathcal{W}_{\rho}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta) \subset \mathcal{T} \mathcal{W}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)
$$

Thus, by Theorems 6 and 7 we have the following results.

Corollary 9 Let the sequence $\left\{d_{n}-(1-\gamma)|\rho|^{n-p}\right\}$, where $d_{n}$ is defined by (11), be positive. The radius of starlikeness of order $\alpha$ for the class $T W_{\rho}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)$ is given by

$$
R_{\alpha}^{*}\left(\mathcal{T} \mathcal{W}_{\rho}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)\right)=\inf _{n \geq k}\left(\frac{(p-\alpha) d_{n}}{(n-\alpha)(1-\gamma)}\right)^{\frac{1}{n-p}}
$$

The functions $f_{n, \eta}$ of the form (32) are the extremal functions.
Corollary 10 Let the sequence $\left\{d_{n}-(1-\gamma)|\rho|^{n-p}\right\}$, where $d_{n}$ is defined by (11), be positive. The radius of convexity of order $\alpha$ for the class $\mathcal{T} \mathcal{W}_{\rho}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)$ is given by

$$
R_{\alpha}^{c}\left(\mathcal{T} \mathcal{W}_{\rho}(p, k ; \phi, \varphi ; \gamma, \delta ; \eta)\right)=\inf _{n \geq k}\left(\frac{(p-\alpha) d_{n}}{n(n-\alpha)(1-\gamma)}\right)^{\frac{1}{n-p}}
$$

where $d_{n}$ is defined by (11).
Remark 3 We conclude this paper by observing that, in view of the definitions of investigated classes which is expressed in terms of the convolution of the functions in (10), involving arbitrary sequences, our main results can lead to several additional new results. In fact, by appropriately selecting these arbitrary sequences, the results presented in this paper would find further applications for the class of analytic functions which would incorporate linear operators. Some of these results were obtained in earlier works, see for example [16, 17, 18, 19, 20].

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[^1]
# Monotonic continuous solution for a mixed type integral inclusion of fractional order 

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Abstract: In this paper we are concerned with the mixed type integral inclusion

$$
x(t) \in p(t)+\int_{0}^{1} k(t, s) F_{1}\left(s, I^{\beta} f_{2}(s, x(s))\right) d s, \quad t \in[0,1] .
$$

The existence of monotonic continuous solution will be proved. As an application the initial-value problem of the arbitrary (fractional) orders differential inclusion

$$
\frac{d x(t)}{d t} \in p(t)+\int_{0}^{1} k(t, s) F_{1}\left(s, D^{\alpha} x(s)\right) d s, \quad \text { a.e., } \quad t>0
$$

will be studied.
AMS Subject Classification:
Key Words and Phrases: Fractional calculus; Caratheodory condition; fixed point theorem; mixed type integral inclusion.

## 1 Introduction

The existence of monotonic integrable solution for the mixed type nonlinear integral equation

$$
\begin{equation*}
x(t)=p(t)+\int_{0}^{1} k(t, s) f_{1}\left(s, I^{\beta} f_{2}(s, x(s)) d s, \quad t \in[0,1], \quad \beta>0\right. \tag{1}
\end{equation*}
$$

has been studied in [6] where the given function $P$ is nondecreasing on $[0,1]$ and the two functions $f_{1}$ and $f_{2}$ are monotonic nondecreasing (in both variables) and satisfy Caratheodory condition.

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Here we relax the condition of monotonicity on the two functions $f_{1}$ and $f_{2}$ and prove the existence of positive continuous solution of (1).
When the given function $p$ is nondecreasing and the kernel $k(t, s)$ is nondecreasing in $t, t \in[0,1]$, we prove that the solution of (1) is nondecreasing.
As a generalization of our results we study the existence of positive monotonic continuous solution of the mixed type integral inclusion

$$
\begin{equation*}
x(t) \in p(t)+\int_{0}^{1} k(t, s) F_{1}\left(s, I^{\beta} f_{2}(s, x(s))\right) d s, \quad t \in[0,1], \quad \beta>0 \tag{2}
\end{equation*}
$$

where the set-valued map $F(t,$.$) is lower semicontinuous from R^{+}$into $R^{+}$and $F(.,$. is measurable.
Finally the differential inclusion of arbitrary (fractional) orders

$$
\begin{equation*}
\frac{d x(t)}{d t} \in p(t)+\int_{0}^{1} k(t, s) F_{1}\left(s, D^{\alpha} x(s)\right) d s, \quad \text { a.e., } \quad t>0 \tag{3}
\end{equation*}
$$

with the initial data

$$
\begin{equation*}
x(0)=x_{\circ} \quad \geq 0 \tag{4}
\end{equation*}
$$

will be studied.

## 2 Preliminaries

Let $L^{1}(I)$ be the class of Lebesgue integrable functions defined on the interval $I=$ $[a, b]$, where $0 \leq a<b<\infty$ and let $\Gamma($.$) be the gamma function.$

Definition 2.1 The fractional integral of the function $f \in L^{1}(I)$ of order $\alpha \in R^{+}$is defined by ([7], [9] and [12])

$$
I_{a}^{\alpha} f(t)=\int_{a}^{t} \frac{(t-s)^{\alpha-1}}{\Gamma(\alpha)} f(s) d s
$$

Definition 2.2 The (Caputo) fractional derivative $D^{\alpha}$ of order $\alpha \in(0,1]$ of the absolutely continuous function $g$ is defined as ([2], [9], [10] and [12])

$$
D_{a}^{\alpha} g(t)=I_{a}^{1-\alpha} \frac{d}{d t} g(t), \quad t \in[a, b] .
$$

Now, we shall state the following theorems which are used in the sequel.

## Theorem 2.1 Schauder's fixed-point Theorem [8]

Let $S$ be a convex subset of a Banach space $B$, let the mapping $T: S \rightarrow S$ be compact and continuous. Then $T$ has at least one fixed-point in $S$.

## Theorem 2.2 Arzela -Ascoli Theorem [4]

Let $E$ be a compact metric space and $C(E)$ be the Banach space of real or complex valued continuous function normed by

$$
\|f\|=\sup _{t \in E}|f(t)|
$$

If $A=\left\{f_{n}\right\}$ is a sequence in $C(E)$ such that $f_{n}$ is uniformly bounded and equicontinuous, then $\overline{\mathrm{A}}$ is compact.

## 3 Main results

Let $C(I), I=[0,1]$ be the class of continuous functions defined on $I$.
In this section we present our main result by proving the existence of monotonic positive solution $x \in C(I)$ for the mixed type integral equation (1).
To facilitate our discussion, let us first state the following assumptions:
(i) $p:[0,1] \rightarrow R^{+}$is continuous. There is a positive constant $p$ such that $|p(t)|<p$.
(ii) $f_{i}:[0,1] \times R^{+} \rightarrow R^{+}, i=1,2$ satisfy caratheodory condition i.e. $f$ is measurable in $t$ for any $x \in R^{+}$and continuous in $x$ for almost all $t \in[0,1]$.
There exist two functions $a_{1}, a_{2} \in L^{1}$ and two positive numbers $b_{1}, b_{2}$ such that

$$
\left|f_{i}(t, x)\right| \leq a_{i}(t)+b_{i}|x|, \quad i=1,2 \quad \forall t \in[0,1] \quad \text { and } \quad x \in R^{+}
$$

(iii) $k:[0,1] \times[0,1] \rightarrow R^{+}$is continuous in $t$ for any $s \in[0,1]$ and measurable in $s$ for any $t \in[0,1]$ such that

$$
\int_{0}^{1} k(t, s)\left(a_{1}(s)+b_{1} I^{\beta} a_{2}(s)\right) d s \leq M_{1} \text { and } \int_{0}^{1} k(t, s) s^{\beta} d s<M_{2}
$$

(iv) $b_{1} b_{2} M_{2}<\Gamma(\beta+1)$.

Remark: It must be noticed that assumption (iii) implies that the two functions

$$
\int_{0}^{1} k(t, s)\left(a_{1}(s)+b_{1} I^{\beta} a_{2}(s)\right) d s \quad \text { and } \int_{0}^{1} k(t, s) s^{\beta} d s
$$

are continuous in $t, \quad t \in[0,1]$.
Now, we are in position to formulate and prove our main result.

Theorem 3.1 Let the assumptions (i)-(iv) be satisfied. Then equation (1) has at least one positive solution $x \in C(I)$.

Proof Define the subset $S$ of $C(I)$ by

$$
S=\{x \in C:|x(t)| \leq r\}, \quad t \in[0,1]
$$

where $r$ is a positive constant. It is clear that $S$ is closed and convex.
Let $T$ be an operator defined by

$$
\begin{equation*}
(T x)(t)=p(t)+\int_{0}^{1} k(t, s) f_{1}\left(s, I^{\beta} f_{2}(s, x(s))\right) d s \quad \forall x \in S \tag{5}
\end{equation*}
$$

Assumption (ii) implies that $T: S \rightarrow C$ is continuous in $x$.
Now for every $x \in S$ we have

$$
\begin{gathered}
|(T x)(t)| \leq|p(t)|+\int_{0}^{1} k(t, s)\left|f_{1}\left(s, I^{\beta} f_{2}(s, x(s))\right)\right| d s \\
\leq p+\int_{0}^{1} k(t, s)\left[a_{1}(s)+b_{1} I^{\beta} f_{2}(s, x(s)) \mid\right] d s \\
\leq p+\int_{0}^{1} k(t, s) a_{1}(s) d s+b_{1} \int_{0}^{1} k(t, s) I^{\beta}\left[a_{2}(s)+b_{2}|x(s)|\right] d s \\
\leq p+\int_{0}^{1} k(t, s)\left[a_{1}(s)+b_{1} I^{\beta} a_{2}(s)\right] d s+b_{1} b_{2} \int_{0}^{1} k(t, s) I^{\beta}|x(s)| d s \\
\leq p+M_{1}+\frac{b_{1} b_{2} r}{\Gamma(\beta+1)} \int_{0}^{1} k(t, s) s^{\beta} d s \\
\leq p+M_{1}+\frac{b_{1} b_{2} M_{2} r}{\Gamma(\beta+1)} .
\end{gathered}
$$

Therefore,

$$
\begin{equation*}
|(T x)(t)| \leq p+M_{1}+\frac{b_{1} b_{2} M_{2} r}{\Gamma(\beta+1)} \tag{6}
\end{equation*}
$$

From the last estimate we deduce that

$$
r=\left(p+M_{1}\right)\left(1-\frac{b_{1} b_{2} M_{2}}{\Gamma(\beta+1)}\right)^{-1}
$$

and $T x \in S$ and hence $T S \subset S$.
Also for $t_{1}, t_{2} \in[0,1]$ such that $t_{1}<t_{2}$, we have

$$
(T x)\left(t_{2}\right)-(T x)\left(t_{1}\right)=p\left(t_{2}\right)-p\left(t_{1}\right)+\int_{0}^{1}\left(k\left(t_{2}, s\right)-k\left(t_{1}, s\right)\right) f_{1}\left(s, I^{\beta} f_{2}(s, x(s))\right) d s
$$

Then

$$
\begin{aligned}
& \left|(T x)\left(t_{2}\right)-(T x)\left(t_{1}\right)\right| \leq\left|p\left(t_{2}\right)-p\left(t_{1}\right)\right|+\int_{0}^{1}\left|k\left(t_{2}, s\right)-k\left(t_{1}, s\right)\right| f_{1}\left(s, I^{\beta} f_{2}(s, x(s))\right) \mid d s \\
& \leq\left|p\left(t_{2}\right)-p\left(t_{1}\right)\right|+\int_{0}^{1}\left|k\left(t_{2}, s\right)-k\left(t_{1}, s\right)\right|\left[a_{1}(s)+b_{1}\left|I^{\beta} f_{2}(s, x(s))\right|\right] d s \\
& \leq\left|p\left(t_{2}\right)-p\left(t_{1}\right)\right|+\int_{0}^{1}\left|k\left(t_{2}, s\right)-k\left(t_{1}, s\right)\right| a_{1}(s) d s \\
& \quad+b_{1} \int_{0}^{1}\left|k\left(t_{2}, s\right)-k\left(t_{1}, s\right)\right|\left|I^{\beta} f_{2}(s, x(s))\right| d s \\
& \leq\left|p\left(t_{2}\right)-p\left(t_{1}\right)\right|+\int_{0}^{1}\left|k\left(t_{2}, s\right)-k\left(t_{1}, s\right)\right| a_{1}(s) d s+b_{1} \int_{0}^{1}\left|k\left(t_{2}, s\right)-k\left(t_{1}, s\right)\right| I^{\beta} a_{2}(s) d s \\
& \quad+b_{1} b_{2} \int_{0}^{1}\left|k\left(t_{2}, s\right)-k\left(t_{1}, s\right)\right| I^{\beta}|x(s)| d s
\end{aligned}
$$

$$
\leq\left|p\left(t_{2}\right)-p\left(t_{1}\right)\right|+\int_{0}^{1}\left|k\left(t_{2}, s\right)-k\left(t_{1}, s\right)\right|\left[a_{1}(s)+b_{1} I^{\beta} a_{2}(s)\right] d s
$$

$$
+b_{1} b_{2} r \int_{0}^{1}\left|k\left(t_{2}, s\right)-k\left(t_{1}, s\right)\right| \int_{0}^{s} \frac{(s-\tau)^{\beta-1}}{\Gamma(\beta)} d \tau d s
$$

$$
\leq\left|p\left(t_{2}\right)-p\left(t_{1}\right)\right|+\int_{0}^{1}\left|k\left(t_{2}, s\right)-k\left(t_{1}, s\right)\right|\left[a_{1}(s)+b_{1} I^{\beta} a_{2}(s)\right] d s
$$

$$
+\frac{b_{1} b_{2} r}{\Gamma(\beta+1)} \int_{0}^{1}\left|k\left(t_{2}, s\right)-k\left(t_{1}, s\right)\right| s^{\beta} d s
$$

From the last inequality, the continuity of the function $p$ and assumption (iii) we deduce the equicontinuity of the functions of $T S$ on $[0,1]$. Then by Arzela-Ascoli Theorem the closure of $T S$ is compact.
Now, all conditions of Schauder's fixed-point Theorem are hold, then $T$ has a fixed point in $S$. Hence there exists at least one positive solution $x \in C(I)$ of (1).

Corollary 3.1 Let the assumption (i)-(iv) are satisfied. If the function $p$ is nondecreasing and $k$ is nondecreasing in $t \in I$, then the solution of (1) is nondecreasing. Proof For $t_{1}, t_{2} \in I$ and $t_{1}<t_{2}$, we have

$$
\begin{aligned}
& x\left(t_{1}\right)=p\left(t_{1}\right)+\int_{0}^{1} k\left(t_{1}, s\right) f_{1}\left(s, I^{\beta} f_{2}(s, x(s))\right) d s \\
& \leq p\left(t_{2}\right)+\int_{0}^{1} k\left(t_{2}, s\right) f_{1}\left(s, I^{\beta} f_{2}(s, x(s))\right) d s=x\left(t_{2}\right) .
\end{aligned}
$$

## 4 Mixed type integral inclusion

Consider now the integral inclusion (2), where $F_{1}:[0,1] \times R^{+} \rightarrow 2^{R^{+}}$has nonempty closed convex values.
As an important consequence of the main result we can present the following:
Theorem 4.1 Let the assumptions of Theorem 3.1 are satisfied and the multifunction $F_{1}$ satisfies the following assumptions:
(1) $F_{1}(t, x)$ are non empty, closed and convex for all $(t, x) \in[0,1] \times R^{+}$,
(2) $F_{1}(t,$.$) is lower semicontinuous from R^{+}$into $R^{+}$,
(3) $F_{1}(.,$.$) is measurable,$
(4) There exist a function $a_{1} \in L^{1}$ and a positive number $b_{1}$ such that

$$
\left|F_{1}(t, x)\right| \leq a_{1}(t)+b_{1}|x| \quad \forall t \in[0,1] .
$$

Then there exists at least one positive solution $x \in C(I)$ of the integral inclusion (2).
Proof By conditions (1) - (4) (see [1], [3], [5] and [11]) we can find a selection function $f_{1}$ (Caratheodory function) $f_{1}:[0,1] \times R^{+} \rightarrow R^{+}$such that $f_{1}(t, x) \in$ $F_{1}(t, x)$ for all $(t, x) \in[0,1] \times R^{+}$, this function satisfies condition (ii) of Theorem 3.1.

Clearly all assumption of Theorem 3.1 are hold, then there exists a continuous positive solution $x \in C(I)$ such that
$x(t)-p(t)=\int_{0}^{1} K(t, s) f_{1}\left(s, I^{\beta} f_{2}(s, x(s))\right) d s \in \int_{0}^{1} K(t, s) F_{1}\left(s, I^{\beta} f_{2}(s, x(s))\right) d s$.

Now, we can easily prove the following Corollary.
Corollary 4.1 Let the assumptions of Theorem 4.1 and the Corollary 3.2 are satisfied, then the solution of (1) is nondecreasing.

## 5 Differential inclusion

Consider now the initial value problem of the differential inclusion (3) with the initial data (4).

Theorem 5.1 Let the assumptions of Theorem 4.1 are satisfied, then the initial value problem (3)-(4) has at least one positive nondecreasing solution $x \in C(I)$.
Proof Let $y(t)=\frac{d x(t)}{d t}$, then equation (3) transformed to the integral inclusion

$$
y(t) \in p(t)+\int_{0}^{1} k(t, s) F_{1}\left(s, I^{1-\alpha} y(s)\right) d s
$$

which by Theorem 4.1 has at least one positive solution $y \in C(I)$.
This implies that the existence of the absolutely continuous solution

$$
x(t)=x_{\circ}+\int_{0}^{t} y(s) d s
$$

which is nondecreasing solution of the initial-value problem (3)-(4).

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# On composition operator in the algebra of functions of two variables with bounded total $\Phi$-variation in Schramm sense 

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#### Abstract

In this paper we consider the so called composition operator being a self-mapping of the Banach algebra of the function of two variables with bounded total $\Phi$-variation in the Schramm sense. The main result of the paper characterizes the composition operator mentioned above which has a generating function being Lipschitzian with respect to the second variable. The basic tool used in our considerations is the concept of the left-left regularization.


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## 1 Introduction

The nonlinear composition operator (which is also known as the superposition operator) is frequently used in many branches of nonlinear analysis and its applications. In order to define such an operator let us assume that $I$ is a real interval (bounded or not) and $f(t, x)=f: I \times \mathbb{R} \rightarrow \mathbb{R}$ is a given function. For an arbitrary function $x(t)=x: I \rightarrow \mathbb{R}$ we may assign the function $F x$ defined as $(F x)(t)=f(t, x(t))$ for $t \in I$. The operator $F$ defined in such a way is called the composition operator generated by the function $f(t, x)$.

One of the basic problems considered in the theory of composition operator can be formulated as follows. Let us assume that $S(I)$ is a set (a space, an algebra, etc.) of some functions acting from $I$ into $\mathbb{R}$. One has to formulate assumptions on the
function $f(t, x)$ guaranteeing that the composition operator $F$ generated by $f(t, x)$ transforms $S(I)$ into itself.
Such a problem was solved in a lot of particular cases. We refer to the monograph [1] for more details concerning that problem.

The second important problem concerning the composition operator depends on the characterization of operators being Lipschitzian in a suitable function space. Such a problem in various situations was studied in a lot of papers (for instance [1], [2], [3], [4], [6], [7], [8], [11]).

In the paper we investigate the problem of characterization of the composition operator being a self-mapping of the Banach algebra of functions of two variables with bounded total variation in the Schramm sense. Namely, we show that such a composition operator is Lipschitzian if and only if it is affine.

The results obtained in the paper generalize those obtained for example in the papers [4], [10].

## 2 Preliminaries

In this section we collect all auxiliary facts which will be needed in the sequel. Let be $\mathbb{R}$ the set of real numbers and $\mathbb{R}_{+}=[0, \infty)$. A function $\varphi: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$is said to be $\varphi$-function if it is continuous on $\mathbb{R}_{+}, \varphi(0)=0, \varphi$ is increasing on $\mathbb{R}_{+}$and $\varphi(t) \rightarrow \infty$ as $t \rightarrow \infty$.
Further, let $\Phi=\left\{\phi_{n}\right\}$ be a sequence of $\varphi$-functions. The sequence $\Phi$ is called the $\Phi$-sequence if $\phi_{n}$ is convex and $\phi_{n+1}(t) \leq \phi_{n}(t)$ for $n=1,2, \ldots$ and for $t \in \mathbb{R}$. Apart from this we assume that these series $\sum \phi_{n}$ diverge for each $t>0$.

Next, let us fix an interval $I=[a, b]$. Assume that $u:[a, b] \rightarrow \mathbb{R}$ is a given function. Let $\phi_{n}$ be a $\Phi$-sequence of functions. If $I_{n}=\left[a_{n}, b_{n}\right]$ is a subinterval of the interval $I$ we write $u\left(I_{n}\right)=u\left(b_{n}\right)-u\left(a_{n}\right)$ (for $n=1,2, \ldots$ ).

We say that the function $u$ has the bounded total $\Phi$-variation in the Schramm sense on the interval $[a, b]$ if

$$
\sum_{n} \phi_{n}\left(\left|u\left(I_{n}\right)\right|\right)<\infty
$$

for each sequence $\left\{I_{n}\right\}$ of closed subintervals of $I$ such that the intersection $I_{i} \cap I_{j}$ is empty or is a singleton for all $i, j=1,2, \ldots, i \neq j$.
We introduced the $\Phi=\left\{\phi_{n, m}\right\}$ bidimensional sequence of increasing convex functions, such that $\phi_{n, m}(0)=0$ and $\phi_{n, m}(t)>0$ for $t>0$ and $n, m=1,2, \ldots$. We shall say that $\Phi$ is a $\Phi^{*}$-sequence if $\phi_{n^{\prime}, m^{\prime}}(t) \leq \phi_{n, m}(t)$ for each $n^{\prime} \leq n, m^{\prime} \leq m$ and $t \in[0, \infty)$. If $\sum_{n=1}^{\infty} \sum_{m=1}^{\infty} \phi_{n, m}$ diverge for $t>0$, we will say that $\Phi$ is a $\Phi$-sequence.

In what follows let us assume that $a=\left(a_{1}, c_{1}\right), b=\left(b_{1}, d_{1}\right)$ are two fixed points in the plane $\mathbb{R}^{2}$. Denote by $I_{a}^{b}$ the rectangle generated by the points $a$ and $b$, i.e. $I_{a}^{b}=\left[a_{1}, b_{1}\right] \times\left[c_{1}, d_{1}\right]$.
Further, let us take two sequences $\left\{I_{n}\right\},\left\{J_{m}\right\}$ of the closed subintervals of the intervals $\left[a_{1}, b_{1}\right]$ and $\left[c_{1}, d_{1}\right]$, respectively. In other words, $I_{n}=\left[a_{n}, b_{n}\right](n=1,2, \ldots), J_{m}=$
$\left[c_{m}, d_{m}\right](m=1,2, \ldots)$. Moreover, let $\Phi=\left\{\phi_{n, m}\right\}$ be a fixed double $\Phi$-sequence and let $u: I_{1} \rightarrow \mathbb{R}$ be a given function. Now, fix $x_{2} \in J_{1}=\left[c_{1}, d_{1}\right]$ and consider the function $u\left(\cdot, x_{2}\right): I_{1} \rightarrow \mathbb{R}$. The quantity $V_{\Phi, I_{1}}^{S}$ defined by the formula

$$
\begin{align*}
V_{\Phi, I_{1}}^{S}(u) & =\sup \sum_{n=1}^{\infty} \phi_{n, m}\left(\left|u\left(I_{n}, x_{2}\right)\right|\right) \\
& =\sup \sum_{n=1}^{\infty} \phi_{n, m}\left(\left|u\left(b_{n}, x_{2}\right)-u\left(a_{n}, x_{2}\right)\right|\right) \tag{1}
\end{align*}
$$

is called the $\Phi$-variation in the Schramm sense of the function $u\left(\cdot, x_{2}\right)$. In the case when $V_{\Phi, I_{1}}^{S}(u)<\infty$ we will say that $u$ has a bounded $\Phi$-variation in the sense of Schramm with respect to the first variable (with the second one fixed).
In the same way we define the $\Phi$-variation of the function $u\left(x_{1}, \cdot\right)$ in the Schramm sense, which will be denoted by $V_{\Phi, J_{1}}^{S}$. If $V_{\Phi, J_{1}}^{S}(u)<\infty$ then $u$ is said to be a function with bounded $\Phi$-variation in Schramm sense with respect to the second variable (with the first one fixed).

Additionally, let us explain that the supremum in formula (1) is taken with respect to all sequences $\left\{I_{n}\right\}$ of the closed subintervals of the interval $I_{1}$. Obviously, in a similar way we understand the supremum in the formula of the quantity $V_{\Phi, J_{1}}^{S}$. Below we provide the definition of the main concept introduced in [4].
Definition 2.1. The quantity $V_{\Phi, I_{a}^{b}}^{S}$ defined by the formula

$$
\begin{aligned}
& V_{\Phi, I_{a}^{b}}^{S}(u)=\sup \sum_{n=1}^{\infty} \sum_{m=1}^{\infty} \phi_{n, m}\left(\left|u\left(I_{n}, J_{m}\right)\right|\right) \\
& \quad=\sup \sum_{n=1}^{\infty} \sum_{m=1}^{\infty} \phi_{n, m}\left(\left|u\left(b_{n}, J_{m}\right)-u\left(a_{n}, J_{m}\right)\right|\right) \\
& \quad=\sup \sum_{n=1}^{\infty} \sum_{m=1}^{\infty} \phi_{n, m}\left(\left|u\left(a_{n}, c_{m}\right)+u\left(b_{n}, d_{m}\right)-u\left(a_{n}, d_{m}\right)-u\left(b_{n}, c_{m}\right)\right|\right)
\end{aligned}
$$

is said to be the bidimensional variation in the sense of Schramm of the function $u$.
Now, let us set the quantity $T V_{\Phi}^{S}$ by putting

$$
\begin{equation*}
T V_{\Phi}^{S}(u)=V_{\Phi, I_{1}}^{S}(u)+V_{\Phi, J_{1}}^{S}(u)+V_{\Phi, I_{a}^{b}}^{S}(u) . \tag{2}
\end{equation*}
$$

This quantity is referred to the total $\Phi$-variation of the function $u$ in the Schramm sense. In the case when $T V_{\Phi}^{S}<\infty$ we say that $u$ is a function with bounded total $\Phi$-variation in Schramm sense.
The set of all functions $u: I_{a}^{b} \rightarrow \mathbb{R}$ having a bounded total $\Phi$-variation will be denoted by $B V_{\Phi}^{S}\left(I_{a}^{b}\right)$.

Next, let us consider the functional $P_{\Phi}$ defined on the set $B V_{\Phi}^{S}\left(I_{a}^{b}\right)$ by the formula

$$
\begin{equation*}
P_{\Phi}(f)=\inf \left\{\epsilon>0: T V_{\Phi}^{S}\left(\frac{f}{\epsilon}\right) \leq 1\right\} \tag{3}
\end{equation*}
$$

The main result proved in [4] asserts that the set $B V_{\Phi}^{S}\left(I_{a}^{b}\right)$ forms a Banach algebra with the norm obtained by the formula

$$
\begin{equation*}
\|f\|_{\Phi}^{S}=|f(a)|+P_{\Phi}(f) \tag{4}
\end{equation*}
$$

Our next result depends on the following lemma.
Lemma 2.1. Let $f \in B V_{\Phi}^{S}\left(I_{a}^{b}\right)$ and $\Phi \in \Phi^{*}$. Then $f$ has the following properties:
(a) If $(t, s),\left(t^{\prime}, s^{\prime}\right) \in I_{a}^{b}$ then $\left|f(t, s)-f\left(t^{\prime}, s^{\prime}\right)\right| \leq 4 \phi_{n, m}^{-1}\left(\frac{1}{2}\right) P_{\Phi}(f)$.
(b) If $P_{\Phi}(f)>0$ then $T V_{\Phi}^{S}\left(f / P_{\Phi}(f)\right) \leq 1$.
(c) Let $r>0$. Then $T V_{\Phi}^{S}(f / r) \leq 1$ if and only if $P_{\Phi}(f) \leq r$.
(d) If $r>0$ and $T V_{\Phi}^{S}\left(f / P_{\Phi}(f)\right)=1$ then $P_{\Phi}(f)=r$.

In what follows let us fix arbitrary $f \in B V_{\Phi}^{S}\left(I_{a}^{b}\right)$. Then, the function $f^{*}: I_{a}^{b} \rightarrow \mathbb{R}$ defined by formula

$$
f^{*}\left(x_{1}, x_{2}\right):= \begin{cases}\lim _{\substack{y_{1} \rightarrow x_{1}-0 \\ y_{2} \rightarrow x_{2}-0}} f\left(y_{1}, y_{2}\right) & \text { if }\left(x_{1}, x_{2}\right) \in\left(a_{1}, b_{1}\right] \times\left(c_{1}, d_{1}\right], \\ \lim _{1} \rightarrow x_{1}-0 \\ y_{2} \rightarrow c_{1}+0 \\ \lim _{1} f\left(y_{1}, y_{2}\right) & \text { if } x_{1} \in\left(a_{1}, b_{1}\right] \text { and } x_{2}=c_{1}, \\ y_{1} \rightarrow a_{1}+0 \\ y_{2} \rightarrow x_{2}-0 \\ \lim _{\substack{ \\y_{1} \rightarrow a_{1}+0 \\ y_{2} \rightarrow c_{1}+0}} f\left(y_{1}, y_{2}\right) & \text { if } x_{1}=a_{1} \text { and } x_{2} \in\left(c_{1}, d_{1}\right], \\ \text { if } x_{1}=a_{1} \text { and } x_{2}=c_{1}\end{cases}
$$

will be called the left-left regularization of the function $f$.
The existence of all one-side limits used above was proved in the book [5].
In the sequel we will denote by $G^{-}\left(I_{a}^{b}\right)$ the class of all left-left regularizations of the function $f \in B V_{\Phi}^{S}\left(I_{a}^{b}\right)$. It can be shown that $G^{-}\left(I_{a}^{b}\right)$ forms a linear space ([9]). Apart from this space $G^{-}\left(I_{a}^{b}\right)$ has the structure of a Banach space with respect to the norm

$$
\|f\|=\sup \left\{|f|:(x, y) \in I_{a}^{b}\right\} .
$$

To present the first result of this paper let us denote by $B V_{\Phi, *}^{S}\left(I_{a}^{b}\right)$ the subspace of the space $B V_{\Phi}^{S}\left(I_{a}^{b}\right)$ containing all functions being left-left continuous on $\left(a_{1}, b_{1}\right] \times$ $\left(c_{1}, d_{1}\right]$.
We have the following result.
Lemma 2.2. If $f \in B V_{\Phi}^{S}\left(I_{a}^{b}\right)$ then $f^{*} \in B V_{\Phi, *}^{S}\left(I_{a}^{b}\right)$.
Proof. First, let us note that according to the definition of the left-left regularization, if $f \in B V_{\Phi}^{S}\left(I_{a}^{b}\right)$ then the function $f^{*}$ is left-left continuous on the set $\left(a_{1}, b_{1}\right] \times\left(c_{1}, d_{1}\right]$. We show that $f^{*} \in B V_{\Phi, *}^{S}\left(I_{a}^{b}\right)$, i.e.

$$
T V_{\Phi}^{S}\left(f^{*}\right)=V_{\Phi, I_{1}}^{S}\left(f^{*}\right)+V_{\Phi, J_{1}}^{S}\left(f^{*}\right)+V_{\Phi, I_{a}^{b}}^{S}\left(f^{*}\right)<\infty
$$

(cf. formulas (1) and (2)).

On composition operator in the algebra of...
At the beginning we show that $V_{\Phi, I_{1}}^{S}\left(f^{*}\right)<\infty$.
To do this we fix $\epsilon>0$ and take a partition $\pi$ of the interval $I_{1}$ generated by the points $a_{1}=t_{0}<t_{1}<\cdots<t_{n}=b_{1}$. Then by virtue of the definition of $f^{*}$ we can find $t_{i}^{\prime} \in\left(t_{i-1}, t_{i}\right) \subset\left[t_{i-1}, t_{i}\right]=\hat{I}_{i}(i=1,2, \ldots, n)$ and $t_{0}^{\prime} \in\left(a_{1}, t_{1}^{\prime}\right), s_{0} \in\left(c_{1}, d_{1}\right) \subset J_{1}$ such that

$$
\left|f^{*}\left(t_{i}, c_{1}\right)-f^{*}\left(t_{i-1}, c_{1}\right)\right| \leq\left|f\left(t_{i}^{\prime}, s_{0}\right)-f\left(t_{i-1}^{\prime}, s_{0}\right)\right|+\frac{1}{4} \phi_{n, m}^{-1}\left(\frac{\epsilon}{m}\right)
$$

Hence, keeping in mind that $\phi_{n, m}$ is increasing, we deduce the following estimate

$$
\begin{align*}
\phi_{n, m} & \left(\left|f^{*}\left(t_{i}, c_{1}\right)-f^{*}\left(t_{i-1}, c_{1}\right)\right|\right) \\
& \leq \phi_{n, m}\left(\left|f\left(t_{i}^{\prime}, s_{0}\right)-f\left(t_{i-1}^{\prime}, s_{0}\right)\right|+\frac{1}{4} \phi_{n, m}^{-1}\left(\frac{\epsilon}{m}\right)\right) \\
& =\phi_{n, m}\left(2\left[\frac{\left|f\left(t_{i}^{\prime}, s_{0}\right)-f\left(t_{i-1}^{\prime}, s_{0}\right)\right|}{2}+\frac{1}{2}\left(\frac{\phi_{n, m}^{-1}\left(\frac{\epsilon}{m}\right)}{4}\right)\right]\right) \\
& \leq \frac{1}{2} \phi_{n, m}\left(2\left|f\left(t_{i}^{\prime}, s_{0}\right)-f\left(t_{i-1}^{\prime}, s_{0}\right)\right|\right)+\frac{1}{2} \phi_{n, m}\left(\frac{1}{2} \phi_{n, m}^{-1}\left(\frac{\epsilon}{m}\right)\right), \tag{5}
\end{align*}
$$

which is also a consequence of the convexity of $\phi_{n, m}$.
On the other hand, since $\phi_{n, m}^{-1}$ is concave, we have

$$
\begin{align*}
\phi_{n, m}^{-1}\left(\frac{\epsilon}{m}\right) & =\phi_{n, m}^{-1}\left(2 \cdot\left(\frac{\epsilon}{2 m}\right)\right) \\
& \leq 2 \phi_{n, m}^{-1}\left(\frac{\epsilon}{2 m}\right) \tag{6}
\end{align*}
$$

From (5) i (6) we obtain

$$
\begin{aligned}
\phi_{n, m} & \left(\left|f^{*}\left(t_{i}, c_{1}\right)-f^{*}\left(t_{i-1}, c_{1}\right)\right|\right) \\
& \leq \frac{1}{2} \phi_{n, m}\left(2\left|f\left(t_{i}^{\prime}, s_{0}\right)-f\left(t_{i-1}^{\prime}, s_{0}\right)\right|\right)+\frac{1}{2} \phi_{n, m}\left(\phi_{n, m}^{-1}\left(\frac{\epsilon}{2 m}\right)\right) \\
& \leq \phi_{n, m}\left(2\left|f\left(t_{i}^{\prime}, s_{0}\right)-f\left(t_{i-1}^{\prime}, s_{0}\right)\right|\right)+\phi_{n, m}\left(\phi_{n, m}^{-1}\left(\frac{\epsilon}{2 m}\right)\right) \\
& =\phi_{n, m}\left(2\left|f\left(t_{i}^{\prime}, s_{0}\right)-f\left(t_{i-1}^{\prime}, s_{0}\right)\right|\right)+\frac{1}{2} \frac{\epsilon}{m}
\end{aligned}
$$

Consequently, we get

$$
\begin{aligned}
\sum_{n=1}^{k} \phi_{n, m}\left(\left|f^{*}\left(\hat{I}_{i}, c_{1}\right)\right|\right) & \leq \sum_{n=1}^{k} \phi_{n, m}\left(2\left|f\left(t_{i}^{\prime}, s_{0}\right)-f\left(t_{i-1}^{\prime}, s_{0}\right)\right|\right)+\epsilon \\
& \leq V_{\Phi, I_{1}}^{S}\left(2 f\left(\cdot, s_{0}\right)\right)+\epsilon \\
& \leq V_{\Phi, I_{1}}^{S}\left(2 f\left(\cdot, d_{1}\right)\right)+\epsilon
\end{aligned}
$$

since $c_{1}<s_{0}<d_{1}$. The last estimate allows us to derive the following one

$$
\begin{equation*}
V_{\Phi, I_{1}}^{S}\left(f^{*}\left(\cdot, c_{1}\right)\right) \leq V_{\Phi, I_{1}}^{S}\left(2 f\left(\cdot, d_{1}\right)\right)<\infty \tag{7}
\end{equation*}
$$

Analogically we can show that

$$
\begin{equation*}
V_{\Phi, J_{1}}^{S}\left(f^{*}\left(a_{1}, \cdot\right)\right) \leq V_{\Phi, J_{1}}^{S}\left(2 f\left(b_{1}, \cdot\right)\right)<\infty . \tag{8}
\end{equation*}
$$

In what follows fix two partitions $\pi_{1}, \pi_{2}$ of the intervals $I_{1}, J_{1}$, respectively, i.e. $\pi_{1}: a_{1}=t_{0}<t_{1}<\cdots<t_{n}=b_{1}, \pi_{2}: c_{1}=s_{0}<s_{1}<\cdots<s_{m}=d_{1}$. In view of the definition of $f^{*}$ we infer that there exist $t_{i}^{\prime} \in\left(t_{i-1}, t_{i}\right) \subset\left[t_{i-1}, t_{i}\right]=\hat{I}_{i}(i=1,2, \ldots, n)$ and $s_{j}^{\prime} \in\left(s_{j-1}, s_{j}\right) \subset\left[s_{j-1}, s_{j}\right]=\hat{J}_{j}(j=1,2, \ldots, m), t_{0}^{\prime} \in\left(a_{1}, t_{1}^{\prime}\right), s_{0}^{\prime} \in\left(c_{1}, s_{1}^{\prime}\right)$ such that

$$
\begin{aligned}
& \quad\left|f^{*}\left(\hat{I}_{i}, \hat{J}_{j}\right)\right|=\left|f^{*}\left(t_{i-1}, s_{j-1}\right)+f^{*}\left(t_{i}, s_{j}\right)-f^{*}\left(t_{i-1}, s_{j}\right)-f^{*}\left(t_{i}, s_{j-1}\right)\right| \\
& \leq\left|f\left(t_{i-1}^{\prime}, s_{j-1}^{\prime}\right)+f\left(t_{i}^{\prime}, s_{j}^{\prime}\right)-f\left(t_{i-1}^{\prime}, s_{j}^{\prime}\right)-f\left(t_{i}^{\prime}, s_{j-1}^{\prime}\right)\right|+\frac{1}{4} \phi_{n, m}^{-1}\left(\frac{\epsilon}{n m}\right) .
\end{aligned}
$$

In a similar way, as earlier, we obtain

$$
\begin{aligned}
& \phi_{n, m}\left(\left|f^{*}\left(\hat{I}_{i}, \hat{J}_{j}\right)\right|\right) \\
& \quad \leq \phi_{n, m}\left(2\left|f\left(t_{i-1}^{\prime}, s_{j-1}^{\prime}\right)+f\left(t_{i}^{\prime}, s_{j}^{\prime}\right)-f\left(t_{i-1}^{\prime}, s_{j}^{\prime}\right)-f\left(t_{i}^{\prime}, s_{j-1}^{\prime}\right)\right|\right)+\frac{\epsilon}{n m}
\end{aligned}
$$

This yields

$$
\begin{array}{r}
\sum_{n=1}^{k} \sum_{m=1}^{l} \phi_{n, m}\left(\left|f^{*}\left(\hat{I}_{i}, \hat{J}_{j}\right)\right|\right) \leq \sum_{n=1}^{k} \sum_{m=1}^{l} \phi_{n, m}\left(2 \mid f\left(t_{i-1}^{\prime}, s_{j-1}^{\prime}\right)+f\left(t_{i}^{\prime}, s_{j}^{\prime}\right)\right. \\
\left.-f\left(t_{i-1}^{\prime}, s_{j}^{\prime}\right)-f\left(t_{i}^{\prime}, s_{j-1}^{\prime}\right) \mid\right)+\epsilon
\end{array}
$$

Consequently, we get

$$
\begin{equation*}
V_{\Phi}^{S}\left(f^{*}, I_{a}^{b}\right) \leq V_{\Phi}^{S}\left(2 f, I_{a}^{b}\right)+\epsilon<\infty . \tag{9}
\end{equation*}
$$

Finally, combining (7), (8) and (9) we derive

$$
T V_{\Phi}^{S}\left(f^{*}\right) \leq T V_{\Phi}^{S}(2 f)<\infty
$$

which means that $f^{*} \in B V_{\Phi, *}^{S}\left(I_{a}^{b}\right)$.
Thus the proof is complete.

## 3 Main result

In this section we prove the main theorem of the paper.
This result characterizes the composition operator acting from the space $B V_{\Phi}^{S}\left(I_{a}^{b}\right)$ into itself which is Lipschitzian.
Theorem 3.1. Let $\Phi$ be convex and let $H: B V_{\Phi}^{S}\left(I_{a}^{b}\right) \rightarrow B V_{\Phi}^{S}\left(I_{a}^{b}\right)$ be a composition operator generated by the function $h: I_{a}^{b} \times \mathbb{R} \rightarrow \mathbb{R}$, i.e.

$$
(H f)(t, s)=h(t, s, f(t, s)), f \in \mathbb{R}^{I_{a}^{b}} \text { for }(t, s) \in I_{a}^{b}
$$

On composition operator in the algebra of...
If the operator $H$ acts from the space $B V_{\Phi}^{S}\left(I_{a}^{b}\right)$ and is Lipschitzian, then

$$
\begin{equation*}
\left|h\left(x, u_{1}\right)-h\left(x, u_{2}\right)\right| \leq \delta\left|u_{1}-u_{2}\right| \tag{10}
\end{equation*}
$$

for all $x \in I_{a}^{b}$ and $u_{1}, u_{2} \in \mathbb{R}$, where $\delta>0$ is a constant.
Moreover, there exist functions $h_{0}, h_{1} \in B V_{\Phi, *}^{S}\left(I_{a}^{b}\right)$ such that

$$
\begin{equation*}
h^{*}(x, u)=h_{0}(x)+h_{1}(x) u \tag{11}
\end{equation*}
$$

for $x \in I_{a}^{b}$ and $u \in \mathbb{R}$. Conversely, if $h_{0}, h_{1} \in B V_{\Phi, *}^{S}\left(I_{a}^{b}\right)$ are functions such that (11) holds, then $H$ acts from the space $B V_{\Phi}^{S}\left(I_{a}^{b}\right)$ into itself and is Lipschitzian.
Proof. Let us fix arbitrarily $\alpha, \beta \in \mathbb{R}, \alpha<\beta$ and define an auxiliary function $\eta_{\alpha, \beta}: \mathbb{R} \rightarrow \mathbb{R}$

$$
\eta_{\alpha, \beta}(t):= \begin{cases}0 & \text { for } t \leq \alpha  \tag{12}\\ \frac{t-\alpha}{\beta-\alpha} & \text { for } \alpha \leq t \leq \beta \\ 1 & \text { for } t \geq \beta\end{cases}
$$

Keeping in mind that the operator $H: B V_{\Phi}^{S}\left(I_{a}^{b}\right) \rightarrow B V_{\Phi}^{S}\left(I_{a}^{b}\right)$ is Lipschitzian, we infer that there exists a constant $\mu>0$ such that $\left\|H f_{1}-H f_{2}\right\|_{\Phi}^{S} \leq \mu\left\|f_{1}-f_{2}\right\|_{\Phi}^{S}$ for any $f_{1}, f_{2} \in B V_{\Phi}^{S}\left(I_{a}^{b}\right)$. The definition of the norm implies

$$
\begin{equation*}
P_{\Phi}\left(H f_{1}-H f_{2}\right) \leq\left\|H f_{1}-H f_{2}\right\|_{\Phi}^{S} \leq \mu\left\|f_{1}-f_{2}\right\|_{\Phi}^{S} \tag{13}
\end{equation*}
$$

In order to simplify the notation let us put $\mathcal{H}=H f_{1}-H f_{2}$. Then, in view of (13) we get

$$
\begin{equation*}
P_{\Phi}(\mathcal{H}) \leq\|\mathcal{H}\|_{\Phi}^{S} \leq \mu\left\|f_{1}-f_{2}\right\|_{\Phi}^{S} \tag{14}
\end{equation*}
$$

If $\left\|f_{1}-f_{2}\right\|_{\Phi}^{S}>0$ then from Lemma 2.1 (c) in (14) we have

$$
T V_{\Phi}^{S}\left(\frac{\mathcal{H}}{\mu\left\|f_{1}-f_{2}\right\|_{\Phi}^{S}}\right) \leq 1
$$

From the definition of $T V_{\Phi}^{S}$, we infer that

$$
\begin{align*}
& \phi_{n, m}\left(\left|\frac{\mathcal{H}\left(\cdot, c_{1}\right)}{\mu\left\|f_{1}-f_{2}\right\|_{\Phi}^{S}}\right|\right) \leq V_{\Phi, I_{1}}^{S}\left(\frac{\mathcal{H}\left(\cdot, x_{2}\right)}{\mu\left\|f_{1}-f_{2}\right\|_{\Phi}^{S}}\right) \leq 1,  \tag{15}\\
& \phi_{n, m}\left(\left|\frac{\mathcal{H}\left(a_{1}, \cdot\right)}{\mu\left\|f_{1}-f_{2}\right\|_{\Phi}^{S}}\right|\right) \leq V_{\Phi, J_{1}}^{S}\left(\frac{\mathcal{H}\left(x_{1}, \cdot\right)}{\mu\left\|f_{1}-f_{2}\right\|_{\Phi}^{S}}\right) \leq 1 \\
& \phi_{n, m}\left(\left|\frac{\mathcal{H}(\cdot, \cdot)}{\mu\left\|f_{1}-f_{2}\right\|_{\Phi}^{S}}\right|\right) \leq \quad V_{\Phi, I_{a}^{b}}^{S}\left(\frac{\mathcal{H}(\cdot, \cdot)}{\mu\left\|f_{1}-f_{2}\right\|_{\Phi}^{S}}\right) \leq 1 .
\end{align*}
$$

Thus, for any $u_{1}, u_{2} \in \mathbb{R}$ and $a=\left(a_{1}, c_{1}\right), b=\left(b_{1}, d_{1}\right), x=\left(x_{1}, x_{2}\right) \in I_{a}^{b}$ we deduce

$$
\begin{align*}
\mid h\left(x_{1},\right. & \left.x_{2}, u_{1}\right)-h\left(x_{1}, x_{2}, u_{2}\right)|=|\mathcal{H}(x)| \\
= & \mid \mathcal{H}\left(x_{1}, c_{1}\right)-\mathcal{H}\left(a_{1}, c_{1}\right)+\mathcal{H}\left(a_{1}, x_{2}\right)-\mathcal{H}\left(a_{1}, c_{1}\right) \\
& +\mathcal{H}\left(a_{1}, c_{1}\right)-\mathcal{H}\left(a_{1}, x_{2}\right)-\mathcal{H}\left(x_{1}, c_{1}\right)+\mathcal{H}(x)+\mathcal{H}\left(a_{1}, c_{1}\right) \mid \\
\leq & \left|\mathcal{H}\left(x_{1}, c_{1}\right)-\mathcal{H}\left(a_{1}, c_{1}\right)\right|+\left|\mathcal{H}\left(a_{1}, x_{2}\right)-\mathcal{H}\left(a_{1}, c_{1}\right)\right| \\
& +\left|\mathcal{H}\left(a_{1}, c_{1}\right)-\mathcal{H}\left(a_{1}, x_{2}\right)-\mathcal{H}\left(x_{1}, c_{1}\right)+\mathcal{H}(x)\right|+\left|\mathcal{H}\left(a_{1}, c_{1}\right)\right| \\
\leq & 3 \phi_{n, m}^{-1}(1) \mu\left\|f_{1}-f_{2}\right\|_{\Phi}^{S}+\left|\mathcal{H}\left(a_{1}, c_{1}\right)\right| . \tag{16}
\end{align*}
$$

To prove the inequality (10) we consider the following cases:

$$
\begin{aligned}
& \text { (i) } a_{1}<x_{1} \leq b_{1} \quad \text { and } \quad c_{1}<x_{2} \leq d_{1}, \\
& \text { (ii) } a_{1}<x_{1} \leq b_{1} \quad \text { and } \quad x_{2}=c_{1}, \\
& \text { (iii) } x_{1}=a_{1} \quad \text { and } \quad c_{1}<x_{2} \leq d_{1}, \\
& \text { (iv) } x_{1}=a_{1} \quad \text { and } \quad x_{2}=c_{1} .
\end{aligned}
$$

Case (i). Consider the functions $f_{1}, f_{2} \in B V_{\Phi}^{S}\left(I_{a}^{b}\right)$ defined by the formulas $f_{\ell}\left(y_{1}, y_{2}\right):=\left(\eta_{a_{1}, x_{1}}\left(y_{1}\right)+\eta_{c_{1}, x_{2}}\left(y_{2}\right)\right) u_{\ell}$ such that $a_{1} \leq y_{1} \leq b_{1}, c_{1} \leq y_{2} \leq d_{1}$ for $\ell=1,2$. Note that $f_{\ell}(a)=f_{\ell}\left(a_{1}, c_{1}\right)=0$ and $f_{\ell}\left(x_{1}, c_{1}\right)=u_{\ell}$ for $\ell=1,2$, and $|\mathcal{H}(a)|=\left|\mathcal{H}\left(a_{1}, c_{1}\right)\right|=0$.
Let $\varepsilon>0$ such that $T V_{\Phi}^{S}\left(\frac{f_{1}-f_{2}}{\varepsilon}\right)=1$. Next, we get

$$
\begin{aligned}
& V_{\Phi, I_{1}}^{S}\left(\frac{\left(f_{1}-f_{2}\right)}{\varepsilon}\left(\cdot, c_{1}\right)\right)=\sup \left\{\sum_{n=1}^{\infty} \phi_{n, m}\left(\left|\frac{\left(f_{1}-f_{2}\right)}{\varepsilon}\left(I_{n}, c_{1}\right)\right|\right):\left\{I_{n}\right\}\right\} \\
&=\sup \left\{\sum_{n=1}^{\infty} \phi_{n, m}\left(\left|\frac{\left(f_{1}-f_{2}\right)}{\varepsilon}\left(b_{n}, c_{1}\right)-\frac{\left(f_{1}-f_{2}\right)}{\varepsilon}\left(a_{n}, c_{1}\right)\right|\right):\left\{I_{n}\right\}\right\} \\
&=\sup \left\{\sum_{n=1}^{\infty} \phi_{n, m}\left(\left|\frac{\left(b_{n}-a_{n}\right)}{\varepsilon\left(x_{1}-a_{1}\right)}\left(u_{1}-u_{2}\right)\right|\right):\left\{I_{n}\right\}\right\} \\
&=\phi_{n, m}\left(\left|\frac{\left(b_{1}-a_{1}\right)}{\varepsilon\left(x_{1}-a_{1}\right)}\left(u_{1}-u_{2}\right)\right|\right) \\
& V_{\Phi, J_{1}}^{S}\left(\frac{\left(f_{1}-f_{2}\right)}{\varepsilon}\left(a_{1}, \cdot\right)\right) \\
& \quad=\sup \left\{\sum_{m=1}^{\infty} \phi_{n, m}\left(\left|\frac{\left(f_{1}-f_{2}\right)}{\varepsilon}\left(a_{1}, J_{m}\right)\right|\right):\left\{J_{m}\right\}\right\}=0 \\
& V_{\Phi, I_{a}^{b}}^{S}\left(\frac{\left(f_{1}-f_{2}\right)}{\varepsilon}(\cdot, \cdot)\right) \\
&=\sup \left\{\sum_{n=1}^{\infty} \sum_{m=1}^{\infty} \phi_{n, m}\left(\left|\frac{\left(f_{1}-f_{2}\right)}{\varepsilon}\left(I_{n}, J_{m}\right)\right|\right):\left\{I_{n}\right\},\left\{J_{m}\right\}\right\}=0
\end{aligned}
$$

Hence

$$
1=T V_{\Phi}^{S}\left(\frac{f_{1}-f_{2}}{\varepsilon}\right)=\phi_{n, m}\left(\left|\frac{\left(b_{1}-a_{1}\right)}{\varepsilon\left(x_{1}-a_{1}\right)}\left(u_{1}-u_{2}\right)\right|\right) .
$$

Moreover, we deduce

$$
\varepsilon=\frac{\left|b_{1}-a_{1}\right|\left|u_{1}-u_{2}\right|}{\phi_{n, m}^{-1}(1)\left|x_{1}-a_{1}\right|}
$$

On composition operator in the algebra of...
By virtue of Lemma 2.1 (d), we choose $P_{\Phi}\left(f_{1}-f_{2}\right)=\varepsilon$ and derive

$$
\begin{align*}
\left\|f_{1}-f_{2}\right\|_{\Phi}^{S} & =\left|\left(f_{1}-f_{2}\right)(a)\right|+P_{\Phi}\left(f_{1}-f_{2}\right) \\
& =0+\frac{\left|b_{1}-a_{1}\right|\left|u_{1}-u_{2}\right|}{\phi_{n, m}^{-1}(1)\left|x_{1}-a_{1}\right|} \tag{17}
\end{align*}
$$

Now, employing (17) in the inequality (16) we obtain (10), i.e.

$$
\begin{aligned}
\left|h\left(x, u_{1}\right)-h\left(x, u_{2}\right)\right| & \leq 3 \phi_{n, m}^{-1}(1) \mu \frac{\left|b_{1}-a_{1}\right|\left|u_{1}-u_{2}\right|}{\phi_{n, m}^{-1}(1)\left|x_{1}-a_{1}\right|}+\left|\mathcal{H}\left(a_{1}, c_{1}\right)\right| \\
& =3 \mu \frac{\left|b_{1}-a_{1}\right|\left|u_{1}-u_{2}\right|}{\left|x_{1}-a_{1}\right|}+0 \\
& =\delta\left|u_{1}-u_{2}\right| \quad\left(\text { where } \delta=3 \mu \frac{\left|b_{1}-a_{1}\right|}{\left|x_{1}-a_{1}\right|}\right)
\end{aligned}
$$

Case (ii). Define the functions

$$
\begin{equation*}
f_{\ell}\left(y_{1}, y_{2}\right)=\left(\eta_{a_{1}, x_{1}}\left(y_{1}\right)\right) u_{\ell} \text { for } \ell=1,2, a_{1} \leq y_{1} \leq b_{1}, c_{1} \leq y_{2} \leq d_{1} . \tag{18}
\end{equation*}
$$

Observe that $f_{\ell}(a)=f_{\ell}\left(a_{1}, c_{1}\right)=\left(\eta_{a_{1}, x_{1}}\left(a_{1}\right)\right) u_{\ell}=0$ for $\ell=1,2$. As in the case (i), we get

$$
\begin{gathered}
V_{\Phi, I_{1}}^{S}\left(\frac{\left(f_{1}-f_{2}\right)}{\varepsilon}\left(\cdot, c_{1}\right)\right)=\sup \left\{\sum_{n=1}^{\infty} \phi_{n, m}\left(\left|\frac{\left(f_{1}-f_{2}\right)}{\varepsilon}\left(I_{n}, c_{1}\right)\right|\right):\left\{I_{n}\right\}\right\} \\
\\
=\phi_{n, m}\left(\left|\frac{\left(b_{1}-a_{1}\right)}{\varepsilon\left(x_{1}-a_{1}\right)}\left(u_{1}-u_{2}\right)\right|\right) \\
V_{\Phi, J_{1}}^{S}\left(\frac{\left(f_{1}-f_{2}\right)}{\varepsilon}\left(a_{1}, \cdot\right)\right)=0=V_{\Phi, I_{a}^{b}}^{S}\left(\frac{\left(f_{1}-f_{2}\right)}{\varepsilon}(\cdot, \cdot)\right)
\end{gathered}
$$

Fix some arbitrary $\varepsilon>0$ such that

$$
1=T V_{\Phi}^{S}\left(\frac{f_{1}-f_{2}}{\varepsilon}\right)=\phi_{n, m}\left(\left|\frac{\left(b_{1}-a_{1}\right)}{\varepsilon\left(x_{1}-a_{1}\right)}\left(u_{1}-u_{2}\right)\right|\right) .
$$

We obtain

$$
\varepsilon=\frac{\left|b_{1}-a_{1}\right|\left|u_{1}-u_{2}\right|}{\phi_{n, m}^{-1}(1)\left|x_{1}-a_{1}\right|}
$$

Taking $P_{\Phi}\left(f_{1}-f_{2}\right)=\varepsilon$ and using the Lemma 2.1 (d) we get

$$
\begin{equation*}
\left\|f_{1}-f_{2}\right\|_{\Phi}^{S}=0+\frac{\left|b_{1}-a_{1}\right|\left|u_{1}-u_{2}\right|}{\phi_{n, m}^{-1}(1)\left|x_{1}-a_{1}\right|} \tag{19}
\end{equation*}
$$

Employing (19) in the inequality (16) we get (10), i.e.

$$
\begin{aligned}
\left|h\left(x, u_{1}\right)-h\left(x, u_{2}\right)\right| & \leq 3 \phi_{n, m}^{-1}(1) \mu \frac{\left|b_{1}-a_{1}\right|\left|u_{1}-u_{2}\right|}{\phi_{n, m}^{-1}(1)\left|x_{1}-a_{1}\right|}+\left|\mathcal{H}\left(a_{1}, c_{1}\right)\right| \\
& =\delta\left|u_{1}-u_{2}\right| \quad\left(\text { where } \delta=3 \mu \frac{\left|b_{1}-a_{1}\right|}{\left|x_{1}-a_{1}\right|}\right)
\end{aligned}
$$

Case (iii). In this case we proceed in the analogous manner as in the case (ii), for which the functions $f_{1}, f_{2} \in B V_{\Phi}^{S}\left(I_{a}^{b}\right)$ are defined by the formulas

$$
f_{\ell}\left(y_{1}, y_{2}\right)=\left(\eta_{c_{1}, x_{2}}\left(y_{2}\right)\right) u_{\ell} \quad \text { for } \quad \ell=1,2, a_{1} \leq y_{1} \leq b_{1}, c_{1} \leq y_{2} \leq d_{1} .
$$

Case (iv). Consider the functions $f_{1}, f_{2} \in B V_{\Phi}^{S}\left(I_{a}^{b}\right)$ defined by

$$
f_{\ell}\left(y_{1}, y_{2}\right)=\left[2-\eta_{a_{1}, b_{1}}\left(y_{1}\right)-\eta_{c_{1}, d_{1}}\left(y_{2}\right)\right] u_{\ell} \quad \text { for } \quad \ell=1,2
$$

such that $a_{1} \leq y_{1} \leq b_{1}$ and $c_{1} \leq y_{2} \leq d_{1}$.
Observe that

$$
\begin{aligned}
f_{1}(a) & =f_{1}\left(a_{1}, c_{1}\right)=\left[2-\eta_{a_{1}, b_{1}}\left(a_{1}\right)-\eta_{c_{1}, d_{1}}\left(c_{1}\right)\right] u_{1}=[2-0-0] u_{1}=2 u_{1}, \\
f_{2}(a) & =f_{2}\left(a_{1}, c_{1}\right)=\left[2-\eta_{a_{1}, b_{1}}\left(a_{1}\right)-\eta_{c_{1}, d_{1}}\left(c_{1}\right)\right] u_{2}=[2-0-0] u_{2}=2 u_{2}, \\
f_{1}(b) & =f_{1}\left(b_{1}, d_{1}\right)=\left[2-\eta_{a_{1}, b_{1}}\left(b_{1}\right)-\eta_{c_{1}, d_{1}}\left(d_{1}\right)\right] u_{1}=[2-1-1] u_{1}=0, \\
f_{2}(b) & =f_{2}\left(b_{1}, d_{1}\right)=\left[2-\eta_{a_{1}, b_{1}}\left(b_{1}\right)-\eta_{c_{1}, d_{1}}\left(d_{1}\right)\right] u_{2}=[2-1-1] u_{2}=0, \\
\mathcal{H}(a) & =h\left(a, u_{1}\right)-h\left(a, u_{2}\right), \\
\mathcal{H}(b) & =\mathcal{H}\left(b_{1}, d_{1}\right)=0,
\end{aligned}
$$

$$
\begin{aligned}
\left(f_{1}-f_{2}\right)\left(a_{1}, J_{m}\right) & =\left[2-\eta_{c_{1}, d_{1}}\left(d_{m}\right)-2+\eta_{c_{1}, d_{1}}\left(c_{m}\right)\right]\left(u_{1}-u_{2}\right) \\
& =\left[-\frac{d_{m}-c_{1}}{d_{1}-c_{1}}+\frac{c_{m}-c_{1}}{d_{1}-c_{1}}\right]\left(u_{1}-u_{2}\right) \\
& =\left[-\frac{d_{m}-c_{m}}{d_{1}-c_{1}}\right]\left(u_{1}-u_{2}\right), \\
\left(f_{1}-f_{2}\right)\left(I_{n}, c_{1}\right) & =\left[-\frac{b_{n}-a_{n}}{b_{1}-a_{1}}\right]\left(u_{1}-u_{2}\right), \\
\left(f_{1}-f_{2}\right)\left(I_{n}, J_{m}\right) & =0 .
\end{aligned}
$$

Hence

$$
\begin{aligned}
V_{\Phi, I_{1}}^{S}\left(\frac{\left(f_{1}-f_{2}\right)}{\varepsilon}\left(I_{n}, c_{1}\right)\right) & =\phi_{n, m}\left(\left|\frac{\left(u_{1}-u_{2}\right)}{\varepsilon}\right|\right) \\
V_{\Phi, J_{1}}^{S}\left(\frac{\left(f_{1}-f_{2}\right)}{\varepsilon}\left(a_{1}, J_{m}\right)\right) & =\phi_{n, m}\left(\left|\frac{\left(u_{1}-u_{2}\right)}{\varepsilon}\right|\right) \\
V_{\Phi, I_{a}^{b}}^{S}\left(\frac{\left(f_{1}-f_{2}\right)}{\varepsilon}\left(I_{n}, J_{m}\right)\right) & =0 .
\end{aligned}
$$

Therefore

$$
T V_{\Phi}^{S}\left(\frac{f_{1}-f_{2}}{\varepsilon}\right)=2 \phi_{n, m}\left(\left|\frac{\left(u_{1}-u_{2}\right)}{\varepsilon}\right|\right)
$$

Taking $\varepsilon>0$ such that

$$
1=T V_{\Phi}^{S}\left(\frac{f_{1}-f_{2}}{\varepsilon}\right)=2 \phi_{n, m}\left(\left|\frac{\left(u_{1}-u_{2}\right)}{\varepsilon}\right|\right)
$$

On composition operator in the algebra of...
we get the following

$$
\varepsilon=\frac{\left|u_{1}-u_{2}\right|}{\phi_{n, m}^{-1}(1 / 2)} .
$$

Now, we select $P_{\Phi}\left(f_{1}-f_{2}\right)=\varepsilon$ and by virtue of the Lemma 2.1 (d) we get the result

$$
\begin{equation*}
\left\|f_{1}-f_{2}\right\|_{\Phi}^{S}=\frac{2 \phi_{n, m}^{-1}(1 / 2)+1}{\phi_{n, m}^{-1}(1 / 2)}\left|u_{1}-u_{2}\right| . \tag{20}
\end{equation*}
$$

In consequence

$$
\begin{align*}
&\left|h\left(a_{1}, c_{1}, u_{1}\right)-h\left(a_{1}, c_{1}, u_{2}\right)\right|=|\mathcal{H}(a)| \\
&= \mid \mathcal{H}\left(b_{1}, c_{1}\right)-\mathcal{H}\left(a_{1}, c_{1}\right)+\mathcal{H}\left(a_{1}, d_{1}\right)-\mathcal{H}\left(a_{1}, c_{1}\right) \\
&+\mathcal{H}\left(a_{1}, c_{1}\right)-\mathcal{H}\left(a_{1}, d_{1}\right)-\mathcal{H}\left(b_{1}, c_{1}\right)+\mathcal{H}\left(b_{1}, d_{1}\right)-\mathcal{H}\left(b_{1}, d_{1}\right) \mid \\
& \leq\left|\mathcal{H}\left(b_{1}, c_{1}\right)-\mathcal{H}\left(a_{1}, c_{1}\right)\right|+\left|\mathcal{H}\left(a_{1}, d_{1}\right)-\mathcal{H}\left(a_{1}, c_{1}\right)\right| \\
&+\left|\mathcal{H}\left(a_{1}, c_{1}\right)-\mathcal{H}\left(a_{1}, d_{1}\right)-\mathcal{H}\left(b_{1}, c_{1}\right)+\mathcal{H}\left(b_{1}, d_{1}\right)\right|+\left|\mathcal{H}\left(b_{1}, d_{1}\right)\right| \\
& \leq 3 \phi_{n, m}^{-1}(1) \mu \| f_{1}-f_{2}| |_{\Phi}^{S}+\left|\mathcal{H}\left(b_{1}, d_{1}\right)\right| \\
&= 3 \phi_{n, m}^{-1}(1) \mu \frac{2 \phi_{n, m}^{-1}(1 / 2)+1}{\phi_{n, m}^{-1}(1 / 2)}\left|u_{1}-u_{2}\right|+|0| \\
&= \delta\left|u_{1}-u_{2}\right| \quad\left(\text { where } \delta=3 \phi_{n, m}^{-1}(1) \mu \frac{2 \phi_{n, m}^{-1}(1 / 2)+1}{\phi_{n, m}^{-1}(1 / 2)}\right) . \tag{21}
\end{align*}
$$

From the foregoing cases we conclude that $h$ is Lipschitzian.
Next, we show the estimation expressed in (11). Let us fix arbitrarily $x_{1} \in\left(a_{1}, b_{1}\right]$, $x_{2} \in\left(c_{1}, d_{1}\right]$ and put $x=\left(x_{1}, x_{2}\right) \in I_{a}^{b}$. For each $k \in \mathbb{N}$ we consider

$$
\begin{aligned}
& a_{1}<\alpha_{1}<\beta_{1}<\alpha_{2}<\beta_{2}<\alpha_{3}<\beta_{3}<\cdots<\alpha_{k}<\beta_{k}<x_{1} \\
& c_{1}<\overline{\alpha_{1}}<\overline{\beta_{1}}<\overline{\alpha_{2}}<\overline{\beta_{2}}<\bar{\alpha}_{3}<\overline{\beta_{3}}<\cdots<\bar{\alpha}_{k}<\bar{\beta}_{k}<x_{2}
\end{aligned}
$$

with $\eta_{k}:\left[a_{1}, b_{1}\right] \rightarrow[0,1]$ and $\bar{\eta}_{k}:\left[c_{1}, d_{1}\right] \rightarrow[0,1]$ two auxiliaries functions defined by the following formulas

$$
\eta_{k}(t):= \begin{cases}0 & \text { for } a_{1} \leq t \leq \alpha_{1}  \tag{22}\\ \eta_{\alpha_{i}, \beta_{i}}(t) & \text { for } \alpha_{i} \leq t \leq \beta_{i}, \quad i=1,2, \ldots, k, \\ 1-\eta_{\beta_{i}, \alpha_{i+1}}(t) & \text { for } \beta_{i} \leq t \leq \alpha_{i+1}, \quad i=1,2, \ldots, k-1, \\ 1 & \text { for } \beta_{k} \leq t \leq b_{1}\end{cases}
$$

and

$$
\bar{\eta}_{k}(s):= \begin{cases}0 & \text { for } \quad c_{1} \leq s \leq \bar{\alpha}_{1},  \tag{23}\\ \eta_{\bar{\alpha}_{i}, \bar{\beta}_{i}}(s) & \text { for } \bar{\alpha}_{i} \leq s \leq \bar{\beta}_{i}, \quad i=1,2, \ldots, k \\ 1-\eta_{\bar{\beta}_{i}, \bar{\alpha}_{i+1}}(s) & \text { for } \quad \bar{\beta}_{i} \leq s \leq \bar{\alpha}_{i+1}, \quad i=1,2, \ldots, k-1, \\ 1 & \text { for } \quad \bar{\beta}_{k} \leq s \leq d_{1} .\end{cases}
$$

For any $u_{1}, u_{2} \in \mathbb{R}$ we define the functions $f_{1}, f_{2}$ by

$$
f_{\ell}\left(y_{1}, y_{2}\right)=\left[-\eta_{k}\left(y_{1}\right)+\bar{\eta}_{k}\left(y_{2}\right)\right] u_{1}+(2-\ell) u_{2}, \quad a_{1} \leq y_{1} \leq b_{1}, \quad c_{1} \leq y_{2} \leq d_{1}
$$

where $\ell=1,2$.
We denote the intervals $I_{\alpha_{k}}$ and $I_{\bar{\alpha}_{k}}$ by $I_{\alpha_{k}}=\left[\alpha_{k}, \beta_{k}\right] \subset\left[a_{1}, b_{1}\right]=I_{1}$ and $I_{\bar{\alpha}_{k}}=$ $\left[\bar{\alpha}_{k}, \bar{\beta}_{k}\right] \subset\left[c_{1}, d_{1}\right]=J_{1}$, then

$$
f_{1}(\cdot, \cdot)-f_{2}(\cdot, \cdot)=u_{2} \quad \text { and } \quad\left\|f_{1}-f_{2}\right\|_{\Phi}^{S}=\left|u_{2}\right| .
$$

Using inequality (15) we have

$$
\begin{gathered}
\sum_{i=1}^{k} \phi_{i, m}\left(\left|\frac{H\left(\beta_{i}, \bar{\beta}_{i}\right)-H\left(\alpha_{i}, \bar{\beta}_{i}\right)}{\mu\left\|f_{1}-f_{2}\right\|_{\Phi}^{S}}\right|\right) \\
\quad \leq \sup \left\{\sum_{k=1}^{\infty} \phi_{k, m}\left(\left|\frac{\mathcal{H}\left(I_{\alpha_{k}}, \bar{\beta}_{k}\right)}{\mu\left\|f_{1}-f_{2}\right\|_{\Phi}^{S}}\right|\right):\left\{I_{\alpha_{k}}\right\}\right\} \\
\quad \leq V_{\Phi, I_{1}}^{S}\left(\frac{\mathcal{H}\left(I_{\alpha_{k}}, \bar{\beta}_{k}\right)}{\mu\left\|f_{1}-f_{2}\right\|_{\Phi}^{S}}\right) \leq 1
\end{gathered}
$$

Thus

$$
\begin{aligned}
\sum_{i=1}^{k} \phi_{i, m}\left(\frac{\mid h\left(\beta_{i}, \bar{\beta}_{i}, f_{1}\left(\beta_{i}, \bar{\beta}_{i}\right)\right)-h\left(\beta_{i}, \bar{\beta}_{i}, f_{2}\left(\beta_{i}, \bar{\beta}_{i}\right)\right)-h\left(\alpha_{i}, \bar{\beta}_{i}, f_{1}\left(\alpha_{i}, \bar{\beta}_{i}\right)\right)}{\mu\left|u_{2}\right|}\right. \\
\left.+\frac{h\left(\alpha_{i}, \bar{\beta}_{i}, f_{2}\left(\alpha_{i}, \bar{\beta}_{i}\right)\right) \mid}{\mu\left|u_{2}\right|}\right) \leq 1
\end{aligned}
$$

Since $f_{1}\left(\beta_{i}, \bar{\beta}_{i}\right)=u_{2}, f_{2}\left(\beta_{i}, \bar{\beta}_{i}\right)=0, f_{1}\left(\alpha_{i}, \bar{\beta}_{i}\right)=u_{1}+u_{2}, f_{2}\left(\alpha_{i}, \bar{\beta}_{i}\right)=u_{1}$, we get from the foregoing estimation

$$
\begin{align*}
\sum_{i=1}^{k} \phi_{i, m}\left(\frac{\mid h\left(\beta_{i}, \bar{\beta}_{i}, u_{2}\right)-h\left(\beta_{i}, \bar{\beta}_{i}, 0\right)-h\left(\alpha_{i}, \bar{\beta}_{i}, u_{1}+u_{2}\right)}{\mu\left|u_{2}\right|}\right. \\
\left.+\frac{h\left(\alpha_{i}, \bar{\beta}_{i}, u_{1}\right) \mid}{\mu\left|u_{2}\right|}\right) \leq 1 \tag{24}
\end{align*}
$$

It is great importance to remark that the constant functions of two variables defined on the rectangle $I_{a}^{b}$ belong to the space $B V_{\Phi}^{S}\left(I_{a}^{b}\right)$ since the composition operator $H$ generated by $h$ acts from $B V_{\Phi}^{S}\left(I_{a}^{b}\right)$ into $B V_{\Phi}^{S}\left(I_{a}^{b}\right)$ and the functions $h(\cdot, u)[x \mapsto$ $h(x, u)$ ] belong to the space $B V_{\Phi}^{S}\left(I_{a}^{b}\right)$ for each $u \in \mathbb{R}$. On the other hand, we know from Lemma 2.2 that the regularization left-left in the first two variables $h^{*}(\cdot, u)$ belongs to the space $B V_{\Phi, *}^{S}\left(I_{a}^{b}\right)$ for all $u \in \mathbb{R}$. If we apply limit in $(24)$ when $\left(\alpha_{1}, \bar{\alpha}_{1}\right) \mapsto$ $\left(x_{1}-0, x_{2}-0\right)$ we obtain

$$
\begin{aligned}
\sum_{i=1}^{k} \phi_{i, m}\left(\frac{\mid h^{*}\left(x_{1}, x_{2}, u_{2}\right)-h^{*}\left(x_{1}, x_{2}, 0\right)-h^{*}\left(x_{1}, x_{2}, u_{1}+u_{2}\right)}{\mu\left|u_{2}\right|}\right. \\
\left.+\frac{h^{*}\left(x_{1}, x_{2}, u_{1}\right) \mid}{\mu\left|u_{2}\right|}\right) \leq 1
\end{aligned}
$$

On composition operator in the algebra of...
Without losing generality we fix $i=n$ for $n=1,2, \ldots, k$

$$
\begin{aligned}
& k \phi_{n, m}\left(\frac{\mid h^{*}\left(x_{1}, x_{2}, u_{2}\right)-h^{*}\left(x_{1}, x_{2}, 0\right)-h^{*}\left(x_{1}, x_{2}, u_{1}+u_{2}\right)}{\mu\left|u_{2}\right|}\right. \\
&\left.+\frac{h^{*}\left(x_{1}, x_{2}, u_{1}\right) \mid}{\mu\left|u_{2}\right|}\right) \leq 1 .
\end{aligned}
$$

Hence

$$
\begin{aligned}
& \phi_{n, m}\left(\frac{\mid h^{*}\left(x_{1}, x_{2}, u_{2}\right)-h^{*}\left(x_{1}, x_{2}, 0\right)-h^{*}\left(x_{1}, x_{2}, u_{1}+u_{2}\right)}{\mu\left|u_{2}\right|}\right. \\
&\left.+\frac{h^{*}\left(x_{1}, x_{2}, u_{1}\right) \mid}{\mu\left|u_{2}\right|}\right) \leq \frac{1}{k}
\end{aligned}
$$

Since $k \in \mathbb{N}$ is arbitrary we derive

$$
\begin{aligned}
& \phi_{n, m}\left(\frac{\mid h^{*}\left(x_{1}, x_{2}, u_{2}\right)-h^{*}\left(x_{1}, x_{2}, 0\right)-h^{*}\left(x_{1}, x_{2}, u_{1}+u_{2}\right)}{\mu\left|u_{2}\right|}\right. \\
&\left.+\frac{h^{*}\left(x_{1}, x_{2}, u_{1}\right) \mid}{\mu\left|u_{2}\right|}\right)=0 .
\end{aligned}
$$

Because $\phi_{n, m}$ is convex for $n, m=1,2, \ldots$ and $\phi(t)=0$ only if $t=0$, then

$$
\frac{\left|h^{*}\left(x_{1}, x_{2}, u_{2}\right)-h^{*}\left(x_{1}, x_{2}, 0\right)-h^{*}\left(x_{1}, x_{2}, u_{1}+u_{2}\right)+h^{*}\left(x_{1}, x_{2}, u_{1}\right)\right|}{\mu\left|u_{2}\right|}=0 .
$$

Therefore

$$
h^{*}\left(x_{1}, x_{2}, u_{2}\right)-h^{*}\left(x_{1}, x_{2}, 0\right)-h^{*}\left(x_{1}, x_{2}, u_{1}+u_{2}\right)+h^{*}\left(x_{1}, x_{2}, u_{1}\right)=0
$$

or equivalently

$$
\begin{equation*}
h^{*}\left(x, u_{1}+u_{2}\right)+h^{*}(x, 0)=h^{*}\left(x, u_{1}\right)+h^{*}\left(x, u_{2}\right) \tag{25}
\end{equation*}
$$

for each $x=\left(x_{1}, x_{2}\right) \in\left(a_{1}, b_{1}\right] \times\left(c_{1}, d_{1}\right]$ and all $u_{1}, u_{2} \in \mathbb{R}$.
Let $x_{1} \in\left(a_{1}, b_{1}\right]$ and $x_{2}=c_{1}$, now we consider the following inequalities

$$
\begin{gathered}
a_{1}<\alpha_{1}<\beta_{1}<\alpha_{2}<\beta_{2}<\alpha_{3}<\beta_{3}<\cdots<\alpha_{k}<\beta_{k}<x_{1}, \\
c_{1}<\overline{\alpha_{1}}<\overline{\beta_{1}}<\overline{\alpha_{2}}<\overline{\beta_{2}}<\bar{\alpha}_{3}<\overline{\beta_{3}}<\cdots<\bar{\alpha}_{k}<\bar{\beta}_{k}<d_{1}, k \in \mathbb{N} .
\end{gathered}
$$

We proceed in the similar way as in the result (24). Taking limit when $\left(\alpha_{1}, \bar{\beta}_{1}\right) \mapsto$ $\left(x_{1}-0, x_{2}+0\right)$ in (24) we obtain (25). The cases $x_{1}=a_{1}$ and $x_{2} \in\left(c_{1}, d_{1}\right]$ or $x_{1}=a_{1}$ and $x_{2}=c_{1}$ are similar.
Thus the equation (25) holds for each $x=\left(x_{1}, x_{2}\right) \in I_{a}^{b}$ and for any $u_{1}, u_{2} \in \mathbb{R}$.
Now, we fix $x=\left(x_{1}, x_{2}\right) \in I_{a}^{b}$ and define the mapping $T_{x}: \mathbb{R} \rightarrow \mathbb{R}$ by

$$
T_{x}(u)=h^{*}(x, u)-h^{*}(x, 0) \quad \forall u \in \mathbb{R} .
$$

Note that we can rewrite the expression in (25) as follows

$$
\begin{equation*}
T_{x}\left(u_{1}+u_{2}\right)=T_{x}\left(u_{1}\right)+T_{x}\left(u_{2}\right) \quad \forall u_{1}, u_{2} \in \mathbb{R} \tag{26}
\end{equation*}
$$

This shows that $T_{x}$ is an additive operator.
For any $u_{1}, u_{2} \in \mathbb{R}$ we have

$$
\begin{aligned}
\left|T_{x}\left(u_{1}\right)-T_{x}\left(u_{2}\right)\right| & =\left|h^{*}\left(x, u_{1}\right)-h^{*}(x, 0)-h^{*}\left(x, u_{2}\right)+h^{*}(x, 0)\right| \\
& =\left|h^{*}\left(x, u_{1}\right)-h^{*}\left(x, u_{2}\right)\right| \\
& \leq \mu\left|u_{1}-u_{2}\right|
\end{aligned}
$$

i.e. $T_{x}(\cdot)$ is Lipschitz-continuous on $\mathbb{R}$. Then exists a mapping $h_{1}: I_{a}^{b} \rightarrow \mathbb{R}$ such that

$$
T_{x}(u)=h_{1}(x) u \quad \forall x \in I_{a}^{b}, \quad \forall u \in \mathbb{R}
$$

Taking $h_{0}(x)=h^{*}(x, 0), x \in I_{a}^{b}$ we derive

$$
h^{*}(x, u)=T_{x}(u)+h^{*}(x, 0)=h_{1}(x) u+h_{0}(x)
$$

Since $h_{0}(\cdot)=h^{*}(\cdot, 0), h_{1}(\cdot)=h^{*}(\cdot, 1)-h^{*}(\cdot, 0)$ and Lemma 2.2 we have that $h_{0}, h_{1} \in$ $B V_{\Phi, *}^{S}\left(I_{a}^{b}\right)$. Thus

$$
h^{*}(x, u)=h_{1}(x) u+h_{0}(x) \quad \forall x \in I_{a}^{b}, \forall u \in \mathbb{R} \text { with } h_{0}, h_{1} \in B V_{\Phi, *}^{S}\left(I_{a}^{b}\right)
$$

Sufficient Condition. Suppose that the composition operator $H$ is given by

$$
(H f)(x)=h_{0}(x)+h_{1}(x) f(x), \quad x \in I_{a}^{b}, \quad f \in B V_{\Phi}^{S}\left(I_{a}^{b}\right)
$$

As $B V_{\Phi}^{S}\left(I_{a}^{b}\right)$ is a Banach algebra, then $H$ maps the space $B V_{\Phi}^{S}\left(I_{a}^{b}\right)$ into itself. Further

$$
\begin{align*}
\left\|H\left(f_{1}\right)-H\left(f_{2}\right)\right\|_{\Phi}^{S} & =\left\|h_{0}+h_{1} f_{1}-h_{0}-h_{1} f_{2}\right\|_{\Phi}^{S} \\
& \leq K\left\|h_{1}\right\|_{\Phi}^{S}\left\|f_{1}-f_{2}\right\|_{\Phi}^{S} \\
& =\lambda\left\|f_{1}-f_{2}\right\|_{\Phi}^{S} \quad\left(\text { where } \lambda=K\left\|h_{1}\right\|_{\Phi}^{S}\right) . \tag{27}
\end{align*}
$$

In consequence $H$ is a Lipschitzian operator.

## Remark 3.1.

1) The Theorem 3.1 is valid for the regularization right-right, left-right and rightleft of $h(\cdot, u) \forall u \in \mathbb{R}$.
2) If $h_{0}, h_{1} \in B V_{\Phi, *}^{S}\left(I_{a}^{b}\right)$ and $\left\|h_{1}\right\|_{\Phi}^{S}<1 / K$, then by Principium of contraction of Banach in combination with (27), exists only one function $f \in B V_{\Phi}^{S}\left(I_{a}^{b}\right)$ such that

$$
f(x)=h_{0}(x)+h_{1}(x) f(x) \quad \forall x \in I_{a}^{b} \subset \mathbb{R}
$$

On composition operator in the algebra of...
The following corollary is the immediate consequence of the Theorem 3.1.
Corollary 3.1. Suppose that $h: I_{a}^{b} \times \mathbb{R} \rightarrow \mathbb{R}$ is such that $h^{*}=h$ in $I_{a}^{b} \times \mathbb{R} \rightarrow \mathbb{R}$, and composition operator $H$ maps space $B V_{\Phi}^{S}\left(I_{a}^{b}\right)$ into itself. Then it is Lipschitzian if and only if there exist functions $h_{0}, h_{1} \in B V_{\Phi, *}^{S}\left(I_{a}^{b}\right)$ such that

$$
h(x, u)=h_{0}(x)+h_{1}(x) u \quad \forall x \in I_{a}^{b}, \quad u \in \mathbb{R} .
$$

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# Interior controllability of the Benjamin-Bona-Mahony equation ${ }^{1}$ 

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ABSTRACT: In this paper we prove the interior approximate controllability of the following Generalized Benjamin-Bona-Mahony type equation (BBM) with homogeneous Dirichlet boundary conditions

$$
\left\{\begin{array}{l}
z_{t}-a \Delta z_{t}-b \Delta z=1_{\omega} u(t, x), \quad t \in(0, \tau), \quad x \in \Omega, \\
z(t, x)=0, \quad t \geq 0, \quad x \in \partial \Omega
\end{array}\right.
$$

where $a \geq 0$ and $b>0$ are constants, $\Omega$ is a domain in $\mathbb{R}^{N}, \omega$ is an open nonempty subset of $\Omega, 1_{\omega}$ denotes the characteristic function of the set $\omega$ and the distributed control $u \in L^{2}\left(0, \tau ; L^{2}(\Omega)\right)$. We prove that for all $\tau>0$ and any nonempty open subset $\omega$ of $\Omega$ the system is approximately controllable on $[0, \tau]$. Moreover, we exhibit a sequence of controls steering the system from an initial state to a final state in a prefixed time. As a consequence of this result we obtain the interior approximate controllability of the heat equation by putting $a=0$ and $b=1$.

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## 1 Introduction.

The original Benjamin-Bona-Mahony equation was proposed in [4] for the case $N=1$ as a model for the propagation of long waves. This equation and related types of pseudo-parabolic equations have been studied by many authors. Results about existence and uniqueness of solutions can be found in [3]; the long time behavior of solutions and the existence of attractors were studied e.g. in [5], [7], [8] and [15], and

[^2]the controllability for the case $N=1$ with control in the boundary has been studied in [13]. Recently the BBM equation with boundary conditions has been studied in [6] and [12].

The interior approximate controllability is a well known, fascinating and important subject in systems theory; there are some works done by [14], [16], [17], [18] and [19]. Particularly, Zuazua in [19] proves the interior approximate controllability of the heat equation

$$
\begin{cases}z_{t}=\Delta z+1_{\omega} u(t, x), & \text { in }(0, \tau) \times \Omega  \tag{1.1}\\ z=0, & \text { on }(0, \tau) \times \partial \Omega \\ z(0, x)=z_{0}(x), & \text { in } \Omega\end{cases}
$$

in two different ways. In the first one, he uses the Hahn-Banach theorem, integration by parts, the adjoint equation, the Carleman estimates and the Holmgren Uniqueness Theorem([11]).

The second method is constructive and uses a variational technique: fix the control time $\tau>0$, the initial and final state $z_{0}=0, z_{1} \in L^{2}(\Omega)$ respectively and $\epsilon>0$; the control steering the initial state $z_{0}$ to a ball of radius $\epsilon>0$ and center $z_{1}$ is given by the point in which the following functional achieves its minimum value

$$
J_{\epsilon}\left(\varphi_{\tau}\right)=\frac{1}{2} \int_{0}^{\tau} \int_{\omega} \varphi^{2} d x d t+\epsilon\left\|\varphi_{\tau}\right\|_{L^{2}(\Omega)}-\int_{\Omega} z_{1} \varphi_{\tau}
$$

where $\varphi$ is the solution of the corresponding adjoint equation with initial data $\varphi_{\tau}$.
In this paper we prove the interior approximate controllability of the following Generalized Benjamin-Bona-Mahony type equation (BBM) with homogeneous Dirichlet boundary conditions

$$
\left\{\begin{array}{l}
z_{t}-a \Delta z_{t}-b \Delta z=1_{\omega} u(t, x), \quad t \in(0, \tau), \quad x \in \Omega  \tag{1.2}\\
z(t, x)=0, \quad t \geq 0, \quad x \in \partial \Omega
\end{array}\right.
$$

where $a \geq 0$ and $b>0$ are constants, $\Omega$ is a domain in $\mathbb{R}^{N}, \omega$ is an open nonempty subset of $\Omega, 1_{\omega}$ denotes the characteristic function of the set $\omega$ and the distributed control $u \in L^{2}\left(0, \tau ; L^{2}(\Omega)\right)$.
The controllability of such systems, with the controls acting on the whole set $\Omega$ was studied in [1]; they considered the approximate controllability of the system

$$
\left\{\begin{array}{l}
z_{t}-a \Delta z_{t}-b \Delta z=b_{1}(x) u_{1}+\ldots+b_{m}(x) u_{m}, \quad t \geq 0, \quad x \in \Omega  \tag{1.3}\\
z(t, x)=0, \quad t \geq 0, \quad x \in \partial \Omega
\end{array}\right.
$$

where $b_{i} \in L^{2}(\Omega ; \mathbb{R})$, the control functions $u_{i} \in L^{2}(0, \tau ; \mathbb{R}) ; i=1,2, \ldots, m$ and $\Omega$ is a bounded domain in $\mathbb{R}^{N}(N \geq 1)$. More precisely, they prove the following result: the system (1.3) is approximately controllable on $[0, \tau], \tau>0$ iff each of the following finite dimensional systems are controllable on $[0, \tau]$

$$
\begin{equation*}
y^{\prime}=-\frac{b \lambda_{j}}{1+a \lambda_{j}} y+B_{j} u, \quad y \in R\left(E_{j}\right), \quad j=1,2, \ldots, \infty \tag{1.4}
\end{equation*}
$$

where

$$
B_{j}: \mathbb{R}^{m} \rightarrow R\left(E_{j}\right), \quad B_{j} U=\sum_{i=1}^{\gamma_{j}} \frac{1}{1+a \lambda_{j}} E_{j} b_{i} U_{i},
$$

$\lambda_{j}{ }^{\prime} s$ are the eigenvalues of $-\Delta$ with Dirichlet boundary condition and $\gamma_{j}$ the corresponding multiplicity, $E_{j}{ }^{\prime} s$ are the projections on the corresponding eigenspaces and $R\left(E_{j}\right)$ denotes the range of $E_{j}$. Since $\operatorname{dim} R\left(E_{j}\right)=\gamma_{j}<\infty$, the controllability of (1.4) is equivalent to the following algebraic condition:

$$
\begin{equation*}
\operatorname{Rank}\left[B_{j}\right]=\gamma_{j}, \quad j=1,2, \ldots, \infty \tag{1.5}
\end{equation*}
$$

In this paper, we are interested in the interior approximate controllability of system (1.2). This is an important problem from the applications point of view, and more general since the control is acting only on a subset $\omega$ of $\Omega$. We prove that for all $\tau>0$ and any nonempty open subset $\omega$ of $\Omega$ the system is approximately controllable on $[0, \tau]$. Moreover, we can exhibit a sequence of controls steering the system from an initial state to a final state in a prefixed time (see Theorem 3.2). As a consequence of this result we obtain the interior approximate controllability of the heat equation (1.1) by putting $a=0$ and $b=1$.

The technique given here is simple and based on the following results:
Theorem 1.1 (see Theorem 1.23 from [2], pg. 20) Suppose $\Omega \subset \mathbb{R}^{n}$ is an open, non-empty and connected set, and $f$ is real analytic function in $\Omega$ with $f=0$ on a non-empty open subset $\omega$ of $\Omega$. Then, $f=0$ in $\Omega$.

Lemma 1.1 (see Lemma 3.14 from [9], pg. 62) Let $\left\{\alpha_{j}\right\}_{j \geq 1}$ and $\left\{\beta_{i, j}: i=1,2, \ldots, m\right\}_{j \geq 1}$ be two sequences of real numbers such that: $\alpha_{1}>\alpha_{2}>\alpha_{3} \cdots$. Then

$$
\sum_{j=1}^{\infty} e^{\alpha_{j} t} \beta_{i, j}=0, \quad \forall t \in[0, \tau], \quad i=1,2, \cdots, m
$$

iff

$$
\beta_{i, j}=0, \quad i=1,2, \cdots, m ; j=1,2, \cdots, \infty .
$$

Theorem 1.2 The eigenfunctions of the operator $-\Delta$ with Dirichlet boundary conditions on $\Omega$ are real analytic functions in $\Omega$.

## 2 Abstract Formulation of the Problem

In this section we choose the space in which this problem will be set as an abstract ordinary differential equation.
Let $Z=L^{2}(\Omega)=L^{2}(\Omega, \mathbb{R})$ and consider the linear unbounded operator $A: D(A) \subset$ $Z \rightarrow Z$ defined by $A \phi=-\Delta \phi$, where

$$
D(A)=H^{2}(\Omega, \mathbb{R}) \cap H_{0}^{1}(\Omega, \mathbb{R}) .
$$

The operator $A$ has the following very well known properties: the spectrum of $A$ consists of eigenvalues

$$
\begin{equation*}
0<\lambda_{1}<\lambda_{2}<\cdots<\lambda_{j}<\cdots \quad \text { with } \quad \lambda_{j} \rightarrow \infty \tag{2.1}
\end{equation*}
$$

each one with finite multiplicity $\gamma_{j}$ equal to the dimension of the corresponding eigenspace. Therefore:
a) there exists a complete orthonormal set $\left\{\phi_{j, k}\right\}$ of eigenvectors of A .
b) for all $z \in D(A)$ we have

$$
\begin{equation*}
A z=\sum_{j=1}^{\infty} \lambda_{j} \sum_{k=1}^{\gamma_{j}}<z, \phi_{j, k}>\phi_{j, k}=\sum_{j=1}^{\infty} \lambda_{j} E_{j} z \tag{2.2}
\end{equation*}
$$

where $<\cdot, \cdot\rangle$ is the inner product in $Z$ and

$$
\begin{equation*}
E_{j} z=\sum_{k=1}^{\gamma_{j}}<z, \phi_{j, k}>\phi_{j, k} \tag{2.3}
\end{equation*}
$$

So, $\left\{E_{j}\right\}$ is a family of complete orthogonal projections in $Z$ and

$$
\begin{equation*}
z=\sum_{j=1}^{\infty} E_{j} z, \quad z \in Z \tag{2.4}
\end{equation*}
$$

c) $-A$ generates the analytic semigroup $\left\{e^{-A t}\right\}$ given by

$$
\begin{equation*}
e^{-A t} z=\sum_{j=1}^{\infty} e^{-\lambda_{j} t} E_{j} z \tag{2.5}
\end{equation*}
$$

Hence, the equation (1.3) can be written as an abstract ordinary differential equation in $Z$ as follows

$$
\begin{equation*}
z^{\prime}+a A z^{\prime}+b A z=1_{\omega} u(t), \quad t \in(0, \tau] \tag{2.6}
\end{equation*}
$$

Since $(I+a A)=a\left(A-\left(-\frac{1}{a}\right) I\right)$ and $-\frac{1}{a} \in \rho(A)(\rho(A)$ is the resolvent set of $A)$, then the operator:

$$
I+a A: D(A) \rightarrow Z
$$

is invertible with bounded inverse

$$
(I+a A)^{-1}: Z \rightarrow D(A)
$$

Therefore, the equation (2.6) also can be written as follows

$$
\begin{equation*}
z^{\prime}+b(I+a A)^{-1} A z=(I+a A)^{-1} 1_{\omega} u(t) \quad t \in(0, \tau) \tag{2.7}
\end{equation*}
$$

Moreover, $(I+a A)$ and $(I+a A)^{-1}$ can be written in terms of the eigenvalues of A:

$$
(I+a A) z=\sum_{j=1}^{\infty}\left(1+a \lambda_{j}\right) E_{j} z
$$

$$
(I+a A)^{-1} z=\sum_{j=1}^{\infty} \frac{1}{1+a \lambda_{j}} E_{j} z
$$

Therefore, if we put $B=(I+a A)^{-1}$, the equation (2.7) can be written as follows

$$
\begin{equation*}
z^{\prime}+b B A z=B B_{\omega} u(t), \quad t \in(0, \tau) \tag{2.8}
\end{equation*}
$$

where $B_{\omega} f=1_{\omega} f$ is a linear a bounded operator from $Z$ to $Z$ and $u \in L^{2}\left(0, \tau ; L^{2}(\Omega)\right)=$ $L^{2}(0, \tau ; Z)$.

Now, we formulate a simple proposition.
Proposition 2.1 The operators $b B A$ and $T(t)=e^{-b B A t}$ are given by the following expression

$$
\begin{gather*}
b B A z=\sum_{j=1}^{\infty} \frac{b \lambda_{j}}{1+a \lambda_{j}} E_{j} z  \tag{2.9}\\
T(t) z=e^{-b B A t} z=\sum_{j=1}^{\infty} e^{\frac{-b \lambda_{j}}{1+a \lambda_{j}} t} E_{j} z \tag{2.10}
\end{gather*}
$$

and

$$
\begin{equation*}
\|T(t)\| \leq e^{-\beta t}, \quad t \geq 0 \tag{2.11}
\end{equation*}
$$

where

$$
\begin{equation*}
\beta=\inf _{j \geq 1}\left\{\frac{b \lambda_{j}}{1+a \lambda_{j}}\right\}=\frac{b \lambda_{1}}{1+a \lambda_{1}} \tag{2.12}
\end{equation*}
$$

With this notation the system (2.8) can be written as follows

$$
\begin{equation*}
z^{\prime}=-\mathcal{A} z+B B_{\omega} u(t), \quad t \in(0, \tau] \tag{2.13}
\end{equation*}
$$

where $\mathcal{A}=b B A$.

## 3 Main Theorem

In this section we shall prove the main result of this paper on the controllability of the linear system (2.13). But first we give the definition of approximate controllability for this system. To this end, for all $z_{0} \in Z$ and a control $u \in L^{2}(0, \tau ; Z)$ the equation (2.13) with $z(0)=z_{0}$ has a unique mild solution given by

$$
\begin{equation*}
z(t)=T(t) z_{0}+\int_{0}^{t} T(t-s) B B_{\omega} u(s) d s, \quad 0 \leq t \leq \tau \tag{3.1}
\end{equation*}
$$

Definition 3.1 We say that (2.13) is approximately controllable in $[0, \tau]$ if for all $z_{0}$, $z_{1} \in Z$ and $\epsilon>0$, there exists a control $u \in L^{2}(0, \tau ; Z)$ such that the solution $z(t)$ given by (3.1) satisfies

$$
\begin{equation*}
\left\|z(\tau)-z_{1}\right\| \leq \epsilon \tag{3.2}
\end{equation*}
$$

Consider the following bounded linear operator:

$$
\begin{equation*}
G: L^{2}(0, \tau ; Z) \rightarrow Z, \quad G u=\int_{0}^{\tau} T(\tau-s) B B_{\omega} u(s) d s \tag{3.3}
\end{equation*}
$$

whose adjoint operator $G^{*}: Z \longrightarrow L^{2}(0, \tau ; Z)$ is given by

$$
\begin{equation*}
\left(G^{*} z\right)(s)=\left(B B_{\omega}\right)^{*} T^{*}(\tau-s) z=B_{\omega}^{*} B^{*} T^{*}(\tau-s) z, \quad \forall s \in[0, \tau], \quad \forall z \in Z \tag{3.4}
\end{equation*}
$$

The following lemma is trivial:
Lema 3.1 The equation (2.13) is approximately controllable on $[0, \tau]$ if, and only if, $\overline{\operatorname{Rang}(G)}=Z$.

The following result is well known from linear operator theory:
Lema 3.2 Let $W$ and $Z$ be Hilbert spaces and $G^{*} \in L(Z, W)$ the adjoint operator of the linear operator $G \in L(W, Z)$. Then

$$
\overline{\operatorname{Rang}(G)}=Z \Longleftrightarrow \operatorname{Ker}\left(G^{*}\right)=\{0\} .
$$

As a consequence of the foregoing Lemma one can prove the following result:
Lema 3.3 Let $W$ and $Z$ be Hilbert spaces and $G^{*} \in L(Z, W)$ the adjoint operator of the linear operator $G \in L(W, Z)$. Then $\overline{\operatorname{Rang}(G)}=Z$ if, and only if, one of the following statements holds:
a) $\operatorname{Ker}\left(G^{*}\right)=\{0\}$.
b) $\left\langle G G^{*} z, z\right\rangle>0, z \neq 0$ in $Z$.
c) $\lim _{\alpha \rightarrow 0^{+}} \alpha\left(\alpha I+G G^{*}\right)^{-1} z=0$.
d) $\sup _{\alpha>0}\left\|\alpha\left(\alpha I+G G^{*}\right)^{-1}\right\| \leq 1$.

The following theorem follows directly from (3.4), lemma 3.1 and lemma 3.3.
Theorem 3.1 (2.13) is approximately controllable on $[0, \tau]$ iff

$$
\begin{equation*}
B_{\omega}^{*} B^{*} T^{*}(t) z=0, \quad \forall t \in[0, \tau], \quad \Rightarrow z=0 . \tag{3.5}
\end{equation*}
$$

Theorem 3.2 (Main Result) For all $\tau>0$ and any open nonempty subset $\omega$ of $\Omega$ the system (2.13) is approximately controllable on $[0, \tau]$. Moreover, a sequence of controls steering the system (2.13) from initial state $z_{0}$ to an $\epsilon$ neighborhood of the final state $z_{1}$ at time $\tau>0$ is given by

$$
u_{\alpha}(t)=B_{\omega}^{*} B^{*} T(\tau-t)\left(\alpha I+G G^{*}\right)^{-1}\left(z_{1}-T(\tau) z_{0}\right),
$$

and the error of this approximation $E_{\alpha}$ is given by

$$
E_{\alpha}=\alpha\left(\alpha I+G G^{*}\right)^{-1}\left(z_{1}-T(\tau) z_{0}\right) .
$$

Proof. We shall apply Theorem 3.1 to prove the controllability of system (2.13). To this end, we observe that

$$
T^{*}(t) z=\sum_{j=1}^{\infty} e^{\frac{-b \lambda_{j}}{1+a \lambda_{j}} t} E_{j} z, \quad B_{\omega}^{*}=B_{\omega}^{*} \quad \text { and } \quad B^{*}=B
$$

Then,

$$
\left(B B_{\omega}\right)^{*} T^{*}(t) z=B_{\omega} B T^{*}(t) z=\sum_{j=1}^{\infty} e^{\frac{-b \lambda_{j}}{1+a \lambda_{j}} t} \frac{1}{1+a \lambda_{j}} B_{\omega} E_{j} z=0, \quad \forall t \in[0, \tau] .
$$

Since $\left\{\frac{-b \lambda_{j}}{1+a \lambda_{j}}: j=1,2, \ldots\right\}$ is a decreasing sequence, then from Lemma 1.1 we obtain that

$$
\left(B_{\omega} E_{j} z\right)(x)=\sum_{k=1}^{\gamma_{j}}<z, \phi_{j, k}>1_{\omega} \phi_{j, k}(x)=0, \quad \forall x \in \Omega, \quad j=1,2, \ldots
$$

i.e.,

$$
\sum_{k=1}^{\gamma_{j}}<z, \phi_{j, k}>\phi_{j, k}(x)=0 \quad \forall x \in \omega, \quad j=1,2, \ldots \ldots
$$

Now, from theorem 1.2 we know that $\phi_{j, k}{ }^{\prime} s$ are analytic functions, which implies the analyticity of $E_{j} z$. Then, from Theorem 1.1 we get that

$$
\sum_{k=1}^{\gamma_{j}}<z, \phi_{j, k}>\phi_{j, k}(x)=0 \quad \forall x \in \Omega, \quad j=1,2, \ldots \ldots
$$

Hence, $E_{j} z=0, \quad j=1,2, \ldots$, which implies that $z=0$.
Now, given the initial and the final states $z_{0}$ and $z_{1}$, we consider the sequence of controls

$$
\begin{aligned}
u_{\alpha}(\cdot) & =B_{\omega}^{*} B^{*} T(\tau-\cdot)\left(\alpha I+G G^{*}\right)^{-1}\left(z_{1}-T(\tau) z_{0}\right) \\
& =G^{*}\left(\alpha I+G G^{*}\right)^{-1}\left(z_{1}-T(\tau) z_{0}\right), \quad \alpha>0
\end{aligned}
$$

Then,

$$
\begin{aligned}
G u_{\alpha} & =G G^{*}\left(\alpha I+G G^{*}\right)^{-1}\left(z_{1}-T(\tau) z_{0}\right) \\
& =\left(\alpha I+G G^{*}-\alpha I\right)\left(\alpha I+G G^{*}\right)^{-1}\left(z_{1}-T(\tau) z_{0}\right) \\
& =z_{1}-T(\tau) z_{0}-\alpha\left(\alpha I+G G^{*}\right)^{-1}\left(z_{1}-T(\tau) z_{0}\right)
\end{aligned}
$$

From part c) of Lemma 3.3 we know that

$$
\lim _{\alpha \rightarrow 0^{+}} \alpha\left(\alpha I+G G^{*}\right)^{-1}\left(z_{1}-T(\tau) z_{0}\right)=0
$$

Therefore,

$$
\lim _{\alpha \rightarrow 0^{+}} G u_{\alpha}=z_{1}-T(\tau) z_{0}
$$

i.e.,

$$
\lim _{\alpha \rightarrow 0^{+}}\left\{T(\tau) z_{0}+\int_{0}^{\tau} T(\tau-s) B B_{\omega} u(s) d s\right\}=z_{1}
$$

This completes the proof of the Theorem.

Corollary 3.1 For all $\tau>0$ and all open nonempty subset $\omega$ of $\Omega$ the heat equation (1.1) is approximately controllable on $[0, \tau]$.

Proof. It is enough to take $a=0$ and $b=1$ in the equation (1.2).

## 4 Final Remarks

The original Benjamin -Bona-Mohany Equation is a non-linear one, here we have proved the approximate controllability of the linear part of this equation, which is the fundamental base for the study of the controllability of the non linear BBM equation. So, our next work is concerned with the controllability of non linear BBM equation

$$
\left\{\begin{array}{l}
z_{t}-a \Delta z_{t}-b \Delta z=1_{\omega} u(t, x)+f(t, z, u(t)), \quad t \in(0, \tau), \quad x \in \Omega  \tag{4.1}\\
z(t, x)=0, \quad t \geq 0, \quad x \in \partial \Omega
\end{array}\right.
$$

where $a \geq 0$ and $b>0$ are constants, $\Omega$ is a domain in $\mathbb{R}^{N}, \omega$ is an open nonempty subset of $\Omega, 1_{\omega}$ denotes the characteristic function of the set $\omega$, the distributed control $u \in L^{2}\left(0, \tau ; L^{2}(\Omega)\right)$ and $f(t, z, u(t))$ is a nonlinear perturbation.

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# The Demyanov metric and measurable multifunctions ${ }^{1}$ 

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Abstract: We present in this paper measurability multifunctions in the family of all convex, bounded sets which need not be closed. The Demyanov metric is discussed.

AMS Subject Classification: 46B20, 46C05, 46E30
Key Words and Phrases: Multifunctions,measurability,convex and bounded sets,the Demyanov metric, the Hausdorff metric

## 1 Introduction and Preliminares

We introduction the following family subsets of $\mathbb{R}^{d}$ :

$$
\mathcal{C}^{d}=\left\{A \in \mathbb{R}^{d}: A \neq \emptyset, \text { convex, bounded }\right\} \quad \mathcal{K}^{d}=\left\{A \in \mathcal{C}^{d}: A \text { is compact }\right\}
$$

For any $A \in \mathbb{R}^{d}, u \in \mathbb{R}^{d}$ we denote

$$
p_{A}(u)=\sup _{a \in A}<a, u>\quad, A(u)=\left\{a \in A:<a, u>=p_{A}(u)\right\}
$$

where $<.>$ is the scalar product, and by recurrence

$$
A\left(u_{1}, \ldots, u_{i}\right)=A\left(u_{1}, \ldots, u_{i-1}\right)\left(u_{i}\right)
$$

By $\mathcal{E}$ denote the set of all orthonormal sequences $\left(e_{1}, \ldots, e_{k}\right), 1 \leq k \leq d$. We shall often use a single letter to denote elements of $\mathcal{E}$, like $E=\left(e_{1}, \ldots, e_{k}\right)$.

Let $A, B \subset \mathbb{R}^{d}$. The Hausdorff distance is defined by

$$
\rho_{H}(A, B)=\max \{e(A, B), e(B, A)\}
$$

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where

$$
e(A, B)=\sup _{a \in A} \operatorname{dist}(\mathrm{a}, \mathrm{~B})=\sup _{\mathrm{a} \in \mathrm{~A}} \inf _{\mathrm{b} \in \mathrm{~B}}\|\mathrm{a}-\mathrm{b}\|
$$

This is a metric in the family of closed sets in $\mathbb{R}^{d}$.
In the [4] we give the following formula for the Demyanov's metric which permitts further extensions to the case of convex but not necessarily closed sets. Let

$$
\mathcal{E}^{k}=\left\{\left(e_{1}, \ldots, e_{j}\right) \in \mathcal{E}: j \geq k\right\}
$$

For $A, B \in \mathcal{K}^{d}$ and $1 \leq k \leq d$

$$
\rho_{D}(A, B)=\sup _{E \in \mathcal{E}^{k}} \rho_{H}(A(E), B(E))
$$

The metric is that of uniform convergence in the set of mappings whose arguments are orthonormal systems of vectors in $\mathbb{R}^{d}$ and values are convex sets orthogonal to these vectors.

## 2 Demyanov's metric in $\mathcal{C}^{d}$ and $\mathcal{X}_{\mathcal{U}}$

Let $\mathcal{E}_{0}=\mathcal{E} \cup\{0\}$. We introduce in $\mathcal{C}^{d}$ the following equivalence relation.
Definition 2.1 $A \equiv B$ iff for every $E \in \mathcal{E}_{0}$ we have

$$
A(E) \neq \emptyset \Leftrightarrow B(E) \neq \emptyset
$$

Remark that the set $\mathcal{U}$ of $E \in \mathcal{E}_{0}$ for which $A(E) \neq \emptyset$ is common for all elements of the same equivalence class-this equivalence class will be then denoted as $\mathcal{K}_{\mathcal{U}}^{d}$.

If a set $\mathcal{U} \subset \mathcal{E}_{0}$ corresponds to some equivalence class then it will be called admissible. Any admissible set $\mathcal{U}$ satisfies the following two conditions:
(i) $0 \in \mathcal{U}$
(ii) $\left(e_{1}, \ldots, e_{k}, e_{k+1}\right) \in \mathcal{U} \Rightarrow\left(e_{1}, \ldots, e_{k}\right) \in \mathcal{U}$

The following example showing that the conditions (i),(ii) are not sufficient for a set $\mathcal{U}$ to be admissible.

Example 2.1 Let $d=2$ and $\mathcal{U}=\{0\} \cup S^{1}$
Let $\mathcal{K}_{\mathcal{U}}^{d} \neq \emptyset$. We remark that the closed ball of $\mathbb{R}^{d}$ is an element of $\mathcal{K}_{\mathcal{U}}^{d}$. Let $A \in \mathcal{K}_{\mathcal{U}}^{d}$ and $e_{1}$ be such that $\bar{A}\left(e_{1}\right)$ is an exposed point in $\bar{A}$-the bar over $A$ denotes the closure. Then $\bar{A}\left(e_{1}\right) \in A$. For any $e_{2} \in S^{1}$ orthogonal to $e_{1}$ we have $A\left(e_{1}\right)=A\left(e_{1}, e_{2}\right) \neq \emptyset$, hence $\left(e_{1}, e_{2}\right) \notin \mathcal{U}$.

Remark that the family off all $A(E), E \in \mathcal{E}_{0}$ consists of all faces of the set $A \in \mathcal{C}^{d}$. The following counterpart of Theorem 2.1.2 [5] is valid for arbitrary $A \in \mathcal{C}^{d}$ (the relint $A$ denote the relative interior of set $A$ ).

Theorem 2.1 If $A(E) \neq A(F)$ then relint $A(E) \cap \operatorname{relint} A(F)=\emptyset$. Morever, the family of sets relint $A(E)$, for all $E \in \mathcal{E}_{0}$ provides a decomposition of $A$.

Fix an admissible family $\mathcal{U}$. Let $E=\left(e_{0}, e_{1}, \ldots, e_{k}\right) \in \mathcal{U}$, where $0 \leq k \leq d$. By $\mathcal{Z}_{E}$ we denote the family of all convex, nonempty, relatively open sets $A$ such that

$$
\forall e_{i} \in E, \forall u, v \in A:<u-v, e_{i}>=0
$$

By $\mathcal{A}_{E}$ denote the element of $\mathcal{Z}_{E}$. We prove the following lemma:
Lemma 2.1 The space $\left(\mathcal{Z}_{E}, \rho_{H}\right)$ is a complete
Proof: Let $E=\left(e_{0}, e_{1}, \ldots, e_{k}\right) \in \mathcal{U}$ where $0 \leq k \leq d$. We consider the Cauchy sequence $\left(\mathcal{A}_{E}^{n}\right)$. Then also $\left(\overline{\mathcal{A}}_{E}^{n}\right)$ is a Cauchy sequence and is element of $\mathcal{K}^{d}$. The space $\left(\mathcal{K}^{d}, \rho_{H}\right)$ is complete. Thus $\lim _{n \rightarrow \infty} \rho_{H}\left(\overline{\mathcal{A}}_{E}^{n}, B\right)=0$, where $B \in \mathcal{K}^{d}$.

If $\operatorname{dim} B<d-k$, then the limit is equal $B$. If $\operatorname{dim} B \geq d-k$ then the limit is equal relintB.

The set of all elements $\mathcal{A}$ of the Cartezian product $\prod_{E \in \mathcal{U}} \mathcal{Z}_{E}$ for which the union $\bigcup_{E \in \mathcal{U}} \mathcal{A}_{E}$ is a bounded subset of $\mathbb{R}^{d}$ will by denoted by $\mathcal{X}_{\mathcal{U}}$.
Example 2.2 Let $d=2$ and $\mathcal{U}=\left\{e_{0},\left(e_{o}, e_{1}\right),\left(e_{0}, e_{1}, e_{2}\right)\right\}$, where $e_{0}=\{0\}, e_{1}=$ $(0,1), e_{2}=(1,0)$. We define $\mathcal{A} \in \mathcal{X}_{\mathcal{U}}$ :

$$
\begin{aligned}
\mathcal{A}_{e_{0}} & =\operatorname{intco}\{(1,1),(-1,1),(-1,-1),(1,-1)\} \\
\mathcal{A}_{\left(e_{0}, e_{1}\right)} & =\text { relintco }\left\{\left(\frac{1}{2}, 1\right),(1,1)\right\} \\
\mathcal{A}_{\left(e_{0}, e_{2}\right)} & =\text { relintco }\left\{\left(1,-\frac{1}{2}\right),\left(1, \frac{1}{2}\right\}\right. \\
\mathcal{A}_{\left(e_{0}, e_{1}, e_{2}\right)} & =\{(1,1)\}
\end{aligned}
$$

The union of all $\mathcal{A}_{E}$ is not convex. Remark that putting

$$
\mathcal{A}_{\left(e_{0}, e_{1}\right)}=\text { relintco }\left\{\left(-\frac{1}{2}, 1\right),\left(\frac{1}{2}, 1\right)\right\}
$$

and

$$
\mathcal{A}_{\left(e_{0}, e_{1}, e_{2}\right)}=\left\{\left(\frac{1}{2}, 1\right)\right\}
$$

we get $\mathcal{A}$ for which the union of values is a convex set belonging to $\mathcal{K}_{\mathcal{U}}$,
In $\mathcal{X}_{\mathcal{U}}$ we introduce the following metric

$$
\rho_{P}(\mathcal{A}, \mathcal{B})=\sup _{E \in \mathcal{U}} \rho_{H}\left(\mathcal{A}_{E}, \mathcal{B}_{E}\right)
$$

The following lemma is a standard result about completeness of the space of bounded maps with the metric uniform convergence.

Lemma 2.2 The metric space $\left(\mathcal{X}_{\mathcal{U}}, \rho_{P}\right)$ is complete

Proof: Let $\mathcal{U}$ the admissible set and $\left(\mathcal{A}^{n}\right)$ the Cauchy sequence in $\mathcal{X}_{\mathcal{U}}$. Using the definition of the metric $\rho_{P}$ we have that for all $E \in \mathcal{U},\left(\mathcal{A}_{E}^{n}\right)$ is the Cauchy sequence in $\mathcal{Z}_{E}$. Hence for all $E \in \mathcal{U}, \lim _{n \rightarrow \infty} \rho_{H}\left(\mathcal{A}_{E}^{n}, \mathcal{A}_{E}^{0}\right)=0$, where $\mathcal{A}_{E}^{0} \in \mathcal{Z}_{E}$. Let $\mathcal{A}^{0}=\bigcup_{E \in \mathcal{U}} \mathcal{A}_{E}^{0}$.

We have that $\mathcal{A}^{0} \in \mathcal{Z}_{\mathcal{U}}$ and $\lim _{n \rightarrow \infty} \rho_{P}\left(\mathcal{A}^{n}, \mathcal{A}^{0}\right)=0$
We remark the following fact. The theorem 2.1 says that for $A \in \mathcal{K}_{\mathcal{U}}^{d}$ the union $\bigcup_{E \in \mathcal{U}} \mathcal{A}_{E}$ is equal $A$, so it is convex.

## 3 Measurable multifunctions in $\mathcal{X}_{\mathcal{U}}$

By $(T, \mathcal{M}, \mu)$ will denote a measurable space, i.e., $T$ is a set, $\mathcal{M}$ is a $\sigma$-field and $\mu$ is a measure such that $\mu(A)<\infty$.

Definition 3.1 $F: T \rightarrow \mathcal{X}_{\mathcal{U}}$ is simple if $F$ takes only finitely many values $A_{1}, \ldots, A_{k}$ such that

$$
\left\{t: F(t)=A_{i}\right\} \in \mathcal{M} \text { for } i=1, \ldots, k
$$

Definition 3.2 A multifunction $F: T \rightarrow \mathcal{X}_{\mathcal{U}}$ is a measurable if there are simple multifunctions $F_{n}: T \rightarrow \mathcal{X}_{\mathcal{U}}$ such that

$$
\lim _{n \rightarrow \infty} \rho_{P}\left(F_{n}(t), F(t)\right)=0 \quad \text { a.e. } t \in T
$$

We can now proof the following result
Theorem 3.1 If $F: T \rightarrow \mathcal{K}^{d}$ is measurable, then for any $A \in \mathcal{K}^{d}$ the set

$$
T_{A}=\{t: F(t) \cap A \neq \emptyset\} \in \mathcal{M}
$$

and the multifunction $F_{A}: T_{A} \rightarrow \mathcal{K}^{d}$ defined by $F_{A}(t)=F(t) \cap A$ is measurable
Proof: Let $A_{n}=\left\{x \in \mathbb{R}^{d}: \operatorname{dist}(x, A) \leq \frac{1}{n}\right\}$ for $n=1,2, \ldots$, and $F_{n}$ be a sequence of simple multifunctions such that

$$
\lim _{n \rightarrow \infty} \rho_{P}\left(F_{n}(t), F(t)\right)=0 \text { a.e. }
$$

We define

$$
T_{n}=\bigcup_{j} \bigcap_{i=j}\left\{t: F_{i}(t) \cap A_{n} \neq \emptyset\right\}
$$

Remark that if $t_{0} \in T_{A}$ then $F\left(t_{0}\right) \cap A \neq \emptyset$, and for all $n$ exist $i_{0}(n)$ such that $F_{i}\left(t_{0}\right) \cap A_{n} \neq \emptyset$ for $i \geq i_{0}(n)$. So $t_{0} \in T_{n}$ for any $n$, hence $T_{A} \subset \bigcap_{n} T_{n}$.

We now prove that $\bigcap_{n} T_{n} \subset T_{A}$.
Let $t_{0} \in \bigcap_{n} T_{n}$, then for any $n$ exist $j_{0}(n)$ such that $F_{i}\left(t_{0}\right) \cap A_{n} \neq \emptyset$ for $i \geq j_{0}(n)$. We assume that $j_{0} \rightarrow \infty$ when $n \rightarrow \infty$. Exist a sequence $x_{n} \in F_{j_{0}(n)}\left(s_{0}\right) \cap A_{n}$ for $n=1,2, \ldots$.

Put $\lim _{n \rightarrow \infty} x_{n}=x_{0}$. We have that $x_{0} \in A$ becouse $\operatorname{dist}\left(x_{n}, A\right) \leq \frac{1}{n}$ and

$$
\operatorname{dist}\left(x_{n}, F\left(t_{0}\right)\right) \leq \rho_{H}\left(F_{j_{0}(n)}\left(t_{0}\right), F\left(t_{0}\right)\right) \leq \rho_{P}\left(F_{j_{0}(n)}\left(t_{0}\right), F\left(t_{0}\right)\right) \rightarrow 0
$$

so $x_{0} \in F\left(t_{0}\right)$, hence $T_{A}=\bigcap_{n} T_{n}$.
The $T_{n}$ is measurable then $T_{A}$ is also.
We remark that $G_{n}(t)=F_{n}(t) \cap A_{n}$ for $t \in T_{A}$ is a sequence of simple multifunctions converges to $F_{A}(t)$ for almost everywhere $t \in T_{A}$.

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[^4]
# Certain differential subordinations using a generalized Sălăgean operator and Ruscheweyh operator 

Alina Alb Lupaş<br>Submitted by: Jan Stankiewicz


#### Abstract

In the present paper we define a new operator using the generalized Sălăgean operator and Ruscheweyh operator. Denote by $D R_{\lambda}^{n}$ the Hadamard product of the generalized Sălăgean operator $D_{\lambda}^{n}$ and Ruscheweyh operator $R^{n}$, given by $D R_{\lambda}^{n}: A \rightarrow A, D R_{\lambda}^{n} f(z)=$ $\left(D_{\lambda}^{n} * R^{n}\right) f(z)$ and $A_{n}=\left\{f \in \mathcal{H}(U): f(z)=z+a_{n+1} z^{n+1}+\ldots, z \in U\right\}$ is the class of normalized analytic functions with $A_{1}=A$. We study some differential subordinations regarding the operator $D R_{\lambda}^{n}$.


AMS Subject Classification: 30C45, 30A20, 34A40
Key Words and Phrases: differential subordination, convex function,best dominant, differential operator, convolution product

## 1 Introduction

Denote by $U$ the unit disc of the complex plane, $U=\{z \in \mathbb{C}:|z|<1\}$ and $\mathcal{H}(U)$ the space of holomorphic functions in $U$.

Let

$$
A_{n}=\left\{f \in \mathcal{H}(U): f(z)=z+a_{n+1} z^{n+1}+\ldots, z \in U\right\}
$$

for $n \in \mathbb{N}$ and $A_{1}=A$.
Denote by

$$
K=\left\{f \in A: \operatorname{Re} \frac{z f^{\prime \prime}(z)}{f^{\prime}(z)}+1>0, z \in U\right\}
$$

the class of normalized convex functions in $U$.
If $f$ and $g$ are analytic functions in $U$, we say that $f$ is subordinate to $g$, written $f \prec g$, if there is a function $w$ analytic in $U$, with $w(0)=0$ and $|w(z)|<1$ for all

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$z \in U$, such that $f(z)=g(w(z))$ for all $z \in U$. If $g$ is univalent, then $f \prec g$ if and only if $f(0)=g(0)$ and $f(U) \subseteq g(U)$.

Let $\psi: \mathbb{C}^{3} \times U \rightarrow \mathbb{C}$ and $h$ be an univalent function in $U$. If $p$ is analytic in $U$ and satisfies the (second-order) differential subordination

$$
\begin{equation*}
\psi\left(p(z), z p^{\prime}(z), z^{2} p^{\prime \prime}(z) ; z\right) \prec h(z), \quad \text { for } \quad z \in U \tag{1.1}
\end{equation*}
$$

then $p$ is called a solution of the differential subordination. The univalent function $q$ is called a dominant of the solutions of the differential subordination, or more simply a dominant, if $p \prec q$ for all $p$ satisfying (1.1).

A dominant $\widetilde{q}$ that satisfies $\widetilde{q} \prec q$ for all dominants $q$ of (1.1) is said to be the best dominant of (1.1). The best dominant is unique up to a rotation of $U$.

Definition 1 (Al Oboudi [2]) For $f \in A, \lambda \geq 0$ and $n \in \mathbb{N}$, the operator $D_{\lambda}^{n}$ is defined by $D_{\lambda}^{n}: A \rightarrow A$,

$$
\begin{aligned}
D_{\lambda}^{0} f(z)= & f(z) \\
D_{\lambda}^{1} f(z)= & (1-\lambda) f(z)+\lambda z f^{\prime}(z)=D_{\lambda} f(z) \\
& \cdots \\
D_{\lambda}^{n} f(z)= & (1-\lambda) D_{\lambda}^{n-1} f(z)+\lambda z\left(D_{\lambda}^{n} f(z)\right)^{\prime}=D_{\lambda}\left(D_{\lambda}^{n-1} f(z)\right), \quad \text { for } \quad z \in U .
\end{aligned}
$$

Remark 1 If $f \in A$ and $f(z)=z+\sum_{j=2}^{\infty} a_{j} z^{j}$, then $D_{\lambda}^{n} f(z)=z+\sum_{j=2}^{\infty}[1+(j-1) \lambda]^{n} a_{j} z^{j}$, for $z \in U$.

Remark 2 For $\lambda=1$ in the above definition we obtain the Sălăgean differential operator [5].

Definition 2 (Ruscheweyh [4]) For $f \in A$ and $n \in \mathbb{N}$, the operator $R^{n}$ is defined by $R^{n}: A \rightarrow A$,

$$
\begin{aligned}
R^{0} f(z)= & f(z) \\
R^{1} f(z)= & z f^{\prime}(z) \\
& \cdots \\
(n+1) R^{n+1} f(z)= & z\left(R^{n} f(z)\right)^{\prime}+n R^{n} f(z), \quad \text { for } z \in U .
\end{aligned}
$$

Remark 3 If $f \in A$ and $f(z)=z+\sum_{j=2}^{\infty} a_{j} z^{j}$, then
$R^{n} f(z)=z+\sum_{j=2}^{\infty} C_{n+j-1}^{n} a_{j} z^{j}$, for $z \in U$.
Lemma 1 (Miller and Mocanu [3]) Let $g$ be a convex function in $U$ and let

$$
h(z)=g(z)+n \alpha z g^{\prime}(z), \quad \text { for } \quad z \in U,
$$

where $\alpha>0$ and $n$ is a positive integer.
If

$$
p(z)=g(0)+p_{n} z^{n}+p_{n+1} z^{n+1}+\ldots, \quad \text { for } \quad z \in U
$$

is holomorphic in $U$ and

$$
p(z)+\alpha z p^{\prime}(z) \prec h(z), \quad \text { for } z \in U
$$

then

$$
p(z) \prec g(z)
$$

and this result is sharp.

## 2 Main Results

Definition 3 Let $\lambda \geq 0$ and $n \in \mathbb{N}$. Denote by $D R_{\lambda}^{n}: A \rightarrow A$ the operator given by the Hadamard product (the convolution product) of the generalized Sălăgean operator $D_{\lambda}^{n}$ and the Ruscheweyh operator $R^{n}$ :

$$
D R_{\lambda}^{n} f(z)=\left(D_{\lambda}^{n} * R^{n}\right) f(z)
$$

for any $z \in U$ and each nonnegative integer $n$.
Remark 4 If $f \in A$ and $f(z)=z+\sum_{j=2}^{\infty} a_{j} z^{j}$, then
$D R_{\lambda}^{n} f(z)=z+\sum_{j=2}^{\infty} C_{n+j-1}^{n}[1+(j-1) \lambda]^{n} a_{j}^{2} z^{j}$, for $z \in U$.
Remark 5 For $\lambda=1$ we obtain the Hadamard product $S R^{n}$ [1] of the Sălăgean operator $S^{n}$ and Ruscheweyh operator $R^{n}$.

Theorem 2 Let $g$ be a convex function such that $g(0)=1$ and let $h$ be the function $h(z)=g(z)+z g^{\prime}(z)$, for $z \in U$. If $\lambda \geq 0, n \in \mathbb{N}, f \in A$ and the differential subordination

$$
\begin{equation*}
\frac{n+1}{\lambda z} D R_{\lambda}^{n+1} f(z)-\frac{n(1-\lambda)}{\lambda z} D R_{\lambda}^{n} f(z)-\left(n-1+\frac{1}{\lambda}\right)\left(D R_{\lambda}^{n} f(z)\right)^{\prime} \prec h(z) \tag{2.2}
\end{equation*}
$$

for $z \in U$, holds, then

$$
\left(D R_{\lambda}^{n} f(z)\right)^{\prime} \prec g(z), \quad \text { for } \quad z \in U
$$

and this result is sharp.
Proof. With notation

$$
p(z)=\left(D R_{\lambda}^{n} f(z)\right)^{\prime}=1+\sum_{j=2}^{\infty} C_{n+j-1}^{n}[1+(j-1) \lambda]^{n} j a_{j}^{2} z^{j-1}
$$

and $p(0)=1$, we obtain for $f(z)=z+\sum_{j=2}^{\infty} a_{j} z^{j}$,

$$
\begin{aligned}
p(z)+ & z p^{\prime}(z)=1+\sum_{j=2}^{\infty} C_{n+j-1}^{n}[1+(j-1) \lambda]^{n} j^{2} a_{j}^{2} z^{j-1} \\
= & \frac{n+1}{\lambda z}\left[z+\sum_{j=2}^{\infty} C_{n+j}^{n+1}[1+(j-1) \lambda]^{n+1} a_{j}^{2} z^{j}\right]+\frac{\lambda-n-1}{\lambda} \\
& -\sum_{j=2}^{\infty} C_{n+j-1}^{n}[1+(j-1) \lambda]^{n} a_{j}^{2} z^{j-1}\left(n-1+\frac{1}{\lambda}\right) j \\
& \quad-\sum_{j=2}^{\infty} C_{n+j-1}^{n}[1+(j-1) \lambda]^{n} a_{j}^{2} z^{j-1} \frac{n(1-\lambda)}{\lambda} \\
= & \frac{n+1}{\lambda z} D R_{\lambda}^{n+1} f(z)-\left(n-1+\frac{1}{\lambda}\right)\left(D R_{\lambda}^{n} f(z)\right)^{\prime}-\frac{n(1-\lambda)}{\lambda z} D R_{\lambda}^{n} f(z) .
\end{aligned}
$$

We have $p(z)+z p^{\prime}(z) \prec h(z)$, for $z \in U$. By using Lemma 1 we obtain $p(z) \prec$ $g(z)$, for $z \in U$, i.e. $\left(D R_{\lambda}^{n} f(z)\right)^{\prime} \prec g(z)$, for $z \in U$ and this result is sharp.

Corollary 3 (see [1]) Let $g$ be a convex function such that $g(0)=1$ and let $h$ be the function $h(z)=g(z)+z g^{\prime}(z)$, for $z \in U$. If $n \in \mathbb{N}, f \in A$ and the differential subordination

$$
\begin{equation*}
\frac{1}{z} S R^{n+1} f(z)+\frac{n}{n+1} z\left(S R^{n} f(z)\right)^{\prime \prime} \prec h(z), \quad \text { for } z \in U \tag{2.3}
\end{equation*}
$$

holds, then

$$
\left(S R^{n} f(z)\right)^{\prime} \prec g(z), \quad \text { for } \quad z \in U
$$

and this result is sharp.
Theorem 4 Let $g$ be a convex function, $g(0)=1$ and let $h$ be the function $h(z)=$ $g(z)+z g^{\prime}(z)$, for $z \in U$. If $n \in \mathbb{N}$ and $f \in A$ verifies the differential subordination

$$
\begin{equation*}
\left(D R_{\lambda}^{n} f(z)\right)^{\prime} \prec h(z), \quad \text { for } \quad z \in U \tag{2.4}
\end{equation*}
$$

then

$$
\frac{D R_{\lambda}^{n} f(z)}{z} \prec g(z), \quad \text { for } \quad z \in U
$$

and this result is sharp.
Proof. For $f \in A$ and $f(z)=z+\sum_{j=2}^{\infty} a_{j} z^{j}$ we have

$$
D R_{\lambda}^{n} f(z)=z+\sum_{j=2}^{\infty} C_{n+j-1}^{n}[1+(j-1) \lambda]^{n} a_{j}^{2} z^{j} \quad \text { for, } z \in U
$$

Consider

$$
\begin{aligned}
p(z) & =\frac{D R_{\lambda}^{n} f(z)}{z}=\frac{z+\sum_{j=2}^{\infty} C_{n+j-1}^{n}[1+(j-1) \lambda]^{n} a_{j}^{2} z^{j}}{z} \\
& =1+\sum_{j=2}^{\infty} C_{n+j-1}^{n}[1+(j-1) \lambda]^{n} a_{j}^{2} z^{j-1}
\end{aligned}
$$

We have $p(z)+z p^{\prime}(z)=\left(D R_{\lambda}^{n} f(z)\right)^{\prime}$, for $z \in U$.
Then $\left(D R_{\lambda}^{n} f(z)\right)^{\prime} \prec h(z)$, for $z \in U$, becomes $p(z)+z p^{\prime}(z) \prec h(z)=g(z)+$ $z g^{\prime}(z)$, for $z \in U$. By using Lemma 1 we obtain $p(z) \prec g(z)$, for $z \in U$, i.e. $\frac{D R_{\lambda}^{n} f(z)}{z} \prec g(z)$, for $z \in U$.

Corollary 5 (see [1]) Let $g$ be a convex function, $g(0)=1$ and let $h$ be the function $h(z)=g(z)+z g^{\prime}(z)$, for $z \in U$. If $n \in \mathbb{N}$ and $f \in A$ verifies the differential subordination

$$
\begin{equation*}
\left(S R^{n} f(z)\right)^{\prime} \prec h(z), \quad \text { for } z \in U \tag{2.5}
\end{equation*}
$$

then

$$
\frac{S R^{n} f(z)}{z} \prec g(z), \quad \text { for } z \in U
$$

and this result is sharp.
Theorem 6 Let $g$ be a convex function such that $g(0)=1$ and let $h$ be the function $h(z)=g(z)+z g^{\prime}(z)$, for $z \in U$. If $n \in \mathbb{N}$ and $f \in A$ verifies the differential subordination

$$
\begin{equation*}
\left(\frac{z D R_{\lambda}^{n+1} f(z)}{D R_{\lambda}^{n} f(z)}\right)^{\prime} \prec h(z), \quad \text { for } z \in U \tag{2.6}
\end{equation*}
$$

then

$$
\frac{D R_{\lambda}^{n+1} f(z)}{D R_{\lambda}^{n} f(z)} \prec g(z), \quad \text { for } \quad z \in U
$$

and this result is sharp.
Proof. For $f \in A$ and $f(z)=z+\sum_{j=2}^{\infty} a_{j} z^{j}$ we have

$$
D R_{\lambda}^{n} f(z)=z+\sum_{j=2}^{\infty} C_{n+j-1}^{n}[1+(j-1) \lambda]^{n} a_{j}^{2} z^{j}, \text { for } z \in U
$$

. Consider

$$
\begin{aligned}
p(z) & =\frac{D R_{\lambda}^{n+1} f(z)}{D R_{\lambda}^{n} f(z)}=\frac{z+\sum_{j=2}^{\infty} C_{n+j}^{n+1}[1+(j-1) \lambda]^{n+1} a_{j}^{2} z^{j}}{z+\sum_{j=2}^{\infty} C_{n+j-1}^{n}[1+(j-1) \lambda]^{n} a_{j}^{2} z^{j}} \\
& =\frac{1+\sum_{j=2}^{\infty} C_{n+j}^{n+1}[1+(j-1) \lambda]^{n+1} a_{j}^{2} z^{j-1}}{1+\sum_{j=2}^{\infty} C_{n+j-1}^{n}[1+(j-1) \lambda]^{n} a_{j}^{2} z^{j-1}} .
\end{aligned}
$$

We have $p^{\prime}(z)=\frac{\left(D R_{\lambda}^{n+1} f(z)\right)^{\prime}}{D R_{\lambda}^{n} f(z)}-p(z) \cdot \frac{\left(D R_{\lambda}^{n} f(z)\right)^{\prime}}{D R_{\lambda}^{n} f(z)}$.
Then $p(z)+z p^{\prime}(z)=\left(\frac{z D R_{\lambda}^{n+1} f(z)}{D R_{\lambda}^{n} f(z)}\right)^{\prime}$.
Relation (2.6) becomes $p(z)+z p^{\prime}(z) \prec h(z)=g(z)+z g^{\prime}(z)$, for $z \in U$, and, by using Lemma 1 we obtain $p(z) \prec g(z)$, for $z \in U$, i.e. $\frac{D R_{\lambda}^{n+1} f(z)}{D R_{\lambda}^{n} f(z)} \prec g(z)$, for $z \in U$.

Corollary 7 (see [1]) Let $g$ be a convex function such that $g(0)=1$ and let $h$ be the function $h(z)=g(z)+z g^{\prime}(z)$, for $z \in U$. If $n \in \mathbb{N}$ and $f \in A$ verifies the differential subordination

$$
\begin{equation*}
\left(\frac{z S R^{n+1} f(z)}{S R^{n} f(z)}\right)^{\prime} \prec h(z), \quad \text { for } z \in U \tag{2.7}
\end{equation*}
$$

then

$$
\frac{S R^{n+1} f(z)}{S R^{n} f(z)} \prec g(z), \quad \text { for } \quad z \in U
$$

and this result is sharp.

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# Multiple hypotheses optimal testing for Markov chains and identification subject to the reliability criterion 

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#### Abstract

The problem of many $(L>2)$ hypotheses testing on distributions of a finite state Markov chain is studied. We apply large deviation techniques (LDT). It is demonstrated that this method of investigation in solving the problem of logarithmically asymptotically optimal (LAO) hypotheses testing is easier, compared with the procedure introduced by Haroutunian. The matrix of exponents $\mathbf{E}=\left\{E_{l \mid m}\right\}, \quad m, l=\{1,2, \ldots, L\}$, of error probabilities of the LAO test $E_{l \mid m}(\phi)=\lim _{N \rightarrow \infty}-\frac{1}{N} \log \alpha_{l \mid m}^{(N)}\left(\phi_{N}\right)$, where $\alpha_{l \mid m}^{(N)}\left(\phi_{N}\right)$ for $l \neq m$ is the probability to accept the hypothesis $l$, when the hypothesis $m$ is true, is determined.

Moreover, the identification of distributions for one object and two independent objects via simple homogeneous stationary Markov chains with finite number of states is discussed in the present paper.


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## 1. Introduction

Applications of information-theoretical methods in mathematical statistics are illustrated in the monographs presented by Kullback [9], Csiszár and Körner [4], Blahut [2], Csiszár and Shields [3], Gutman [6] and others. Numerous papers have been devoted to the study of exponential decrease, as the sample size $N$ goes to infinity, of the error probabilities $\alpha_{1}^{(N)}$ of the first kind and $\alpha_{2}^{(N)}$ of the second kind of the optimal tests for two simple statistical hypotheses. Similar problems for Markov dependence
of experiments were investigated by Natarajan [10], Haroutunian [7], [8], Dembo and Zeitouni [5], and others. In the book of Csiszár and Shields [3] different asymptotic aspects of two hypotheses testing for independent identically distributed observations are considered via theory of large deviations.

In this paper, we aim to solve the problem in order to describe the matrix of exponents $E=\left\{E_{l \mid m}\right\}, m, l=\{1,2, \ldots, L\}$ of probabilities $\alpha_{l \mid m}^{(N)}=\exp \left(-N E_{l \mid m}\right)$, where $\alpha_{l \mid m}^{(N)}$ for $l \neq m$ is the probability to accept hypothesis $l$, when hypothesis $m$ is true, for finite state Markov chain by application of large deviation techniques (LDT). We will demonstrated that the solution of the mentioned problem is more concise and hence easier than the procedure introduced by Haroutunian [7].

Ahlswede and Haroutunian [1], formulated an ensemble of problems on multiple hypotheses testing for multiple objects and on identification of hypotheses under reliability requirement. In this paper, we also solve this problem through identification of distributions of many hypotheses for one object and two independent objects, using simple homogeneous stationary finite states of Markov chains.

In Section 2, we present a Theorem of LDT for Markov chains and the result for hypotheses testing and in Section 3, the problem of identification for Markov chain and finally in Section 4, we discuss the general case of the problem of identification of distributions for two independent Markov chains.

## 2. Application of LDT On Many Hypotheses Optimal Testing for Markov chains

Let $\mathbf{x}=\left(x_{0}, x_{1}, x_{2}, \ldots, x_{N}\right), x_{n} \in \mathcal{X}=\{1,2, \ldots, I\}, \mathbf{x} \in \mathcal{X}^{N+1}, N=0,1,2, \ldots$, be vectors of observations of a simple homogeneous stationary Markov chain with finite number $I$ of states. The hypotheses concern the irreducible matrices of the transition probabilities

$$
P_{l}=\left\{P_{l}(j \mid i), i, j=\{1,2, \ldots, I\}\right\}, l=\{1,2, \ldots, L\} .
$$

The stationarity of the chain provides existence for each $l=\{1,2, \ldots, L\}$ of the unique stationary distribution $Q_{l}=\left\{Q_{l}(i), i=\{1,2, \ldots, I\}\right.$, such that

$$
\sum_{i} Q_{l}(i) P_{l}(j \mid i)=Q_{l}(j), \quad \sum_{i} Q_{l}(i)=1, \quad i, j=\{1,2, \ldots, I\}
$$

We define the joint distributions

$$
Q_{l} \circ P_{l}=\left\{Q_{l}(i) P_{l}(j \mid i), i, j=\{1,2, \ldots, I\}\right\}, l=\{1,2, \ldots, L\} .
$$

Let us denote $D\left(Q \circ P \| Q_{l} \circ P_{l}\right)$ Kullback-Leibler divergence $(P$ is an irreducible matrix of transition probabilities of some sttionary Markov chain and $Q$ be the corresponding stationary $P D$ )

$$
D\left(Q \circ P \| Q_{l} \circ P_{l}\right)=\sum_{i, j} Q(i) P(j \mid i)\left[\log Q(i) P(j \mid i)-\log Q_{l}(i) P_{l}(j \mid i)\right]
$$

$$
=D\left(Q \| Q_{l}\right)+D\left(Q \circ P \| Q \circ P_{l}\right),
$$

of distribution

$$
Q \circ P=\{Q(i) P(j \mid i), i, j=\{1,2, \ldots, I\}\}
$$

with respect to distribution $Q_{l} \circ P_{l}$ where

$$
D\left(Q \| Q_{l}\right)=\sum_{i} Q(i)\left[\log Q(i)-\log Q_{l}(i)\right], l=\{1,2, \ldots, L\}
$$

Let us name the second order type of vector $\mathbf{x}$ the square matrix of $I^{2}$ relative frequencies $\left\{N(i, j) N^{-1}, i, j=\{1,2, \ldots, I\}\right\}$ of the simultaneous appearance in $\mathbf{x}$ of the states $i$ and $j$ on the pairs of neighbor places. It is clear that $\sum_{i j} N(i, j)=N$. Denoted by $\mathcal{T}_{Q \circ P}^{N}$, the set of vectors from $\mathcal{X}^{N+1}$ have the second order type in a way that for some joint PD $Q \circ P$

$$
N(i, j)=N Q(i) P(j \mid i), \quad i=\{1,2, \ldots, I\}, \quad j=\{1,2, \ldots, I\} .
$$

The set of all joint $\mathrm{PD} Q \circ P$ on $\mathcal{X}$ is denoted by $\mathcal{Q} \circ \mathcal{P}(\mathcal{X})$ and the set of all possible second order types for joint PD $Q \circ P$ is denoted by $\mathcal{Q} \circ \mathcal{P}^{N}(\mathcal{X})$. Note that if vector $\mathrm{x} \in \mathcal{T}_{Q \circ P}^{N}$, then

$$
\sum_{j} N(i, j)=N Q(i), i=\{1,2, \ldots, I\}, \quad \sum_{i} N(i, j)=N Q^{\prime}(j), j=\{1,2, \ldots, I\}
$$

for somewhat different from $Q \mathrm{PD} Q^{\prime}$, which in accordance with the definition of $N(i, j)$, are closed enough

$$
\left|N Q(i)-N Q^{\prime}(i)\right| \leq 1, i=\{1,2, \ldots, I\},
$$

and in the limit, when $N \rightarrow \infty$, the distribution $Q$ coincides with $Q^{\prime}$ and may be taken as stationary for conditional PD $P$ :

$$
\sum_{i} Q(i) P(j \mid i)=Q(j), j \in \mathcal{X}
$$

The probability of vector $\mathbf{x} \in \mathcal{X}^{N+1}$ of the Markov chain with transition probabilities $P_{l}$ and stationary distribution $Q_{l}$, is the following

$$
\begin{gathered}
Q_{l} \circ P_{l}^{N}(\mathbf{x})=Q_{l}\left(x_{0}\right) \prod_{n=1}^{N} P_{l}\left(x_{n} \mid x_{n-1}\right), l=\{1,2, \ldots, I\} \\
Q_{l} \circ P_{l}^{N}(\mathcal{A})=\bigcup_{\mathbf{x} \in \mathcal{A}} Q_{l} \circ P_{l}^{N}(\mathbf{x}), \mathcal{A} \subset \mathcal{X}^{N+1}
\end{gathered}
$$

Note that for $l=\{1,2, \ldots, L\}$ the probability of $\mathbf{x}$ from $\mathcal{T}_{Q \circ P}^{N}$ can be written as

$$
Q_{l} \circ P_{l}^{N}(\mathbf{x})=Q_{l}\left(x_{0}\right) \prod_{i, j} P_{l}(j \mid i)^{N Q(i) P(j \mid i)}
$$

Note also that if $Q \circ P$ is absolutely continuous relative to $Q_{l} \circ P_{l}$, then

$$
Q_{l} \circ P_{l}^{N}\left(\mathcal{T}_{Q \circ P}^{N}\right)=\exp \left\{-N\left(D\left(Q \circ P \| Q \circ P_{l}\right)\right)+o(1)\right\}
$$

where

$$
\begin{gathered}
o(1)=\max \left(\max _{i}\left|N^{-1} \log Q_{l}(i)\right|: Q_{l}(i)>0\right) \\
\left(\max _{i}\left|N^{-1} \log Q_{l}(i)\right|: Q_{l}(i)>0\right) \rightarrow 0, \quad \text { when } N \rightarrow \infty
\end{gathered}
$$

Indeed, this is not difficult to verify, taking into account the number $\left|\mathcal{T}_{Q \circ P}^{N}\right|$ of vectors in $\mathcal{T}_{Q \circ P}^{N}$ which is equal to

$$
\exp \left\{-N\left(\sum_{i, j} Q(i) P(j \mid i) \log P(j \mid i)\right)+o(1)\right\}
$$

Non-randomized test $\phi_{N}(\mathbf{x})$ accepts one of the hypotheses $H_{l}, l=\{1,2, \ldots, L\}$ on the basis of the trajectory $\mathbf{x}=\left(x_{0}, x_{1}, \ldots, x_{N}\right)$ of the $N+1$ observations. Let us denote $\alpha_{l \mid m}^{(N)}\left(\phi_{N}\right)$ the probability to accept the hypothesis $H_{l}$ under the condition that $H_{m}, m \neq l$, is true. For $l=m$ we denote $\alpha_{m \mid m}^{(N)}\left(\phi_{N}\right)$ the probability to reject the hypothesis $H_{m}$. It is clear that

$$
\begin{equation*}
\alpha_{m \mid m}^{(N)}\left(\phi_{N}\right)=\sum_{l \neq m} \alpha_{l \mid m}^{(N)}\left(\phi_{N}\right), m=\{1,2, \ldots, L\} \tag{1}
\end{equation*}
$$

This probability is called the error probability of the $m$-th kind of the test $\phi_{N}$. The quadratic matrix of $L^{2}$ error probabilities $\left\{\alpha_{l \mid m}^{(N)}(\phi), m, l=\{1,2, \ldots, L\}\right.$ is occasionally called the power of the tests. To every trajectory $\mathbf{x}$, the test $\phi_{N}$ puts in one correspondence from $L$ hypotheses. Thus, the space $\mathcal{X}^{N+1}$ will be divided into $L$ parts,

$$
\mathcal{G}_{l}^{N}=\left\{\mathbf{x}, \phi_{N}(\mathbf{x})=l\right\}, l=\{1,2, \ldots, L\}
$$

and

$$
\alpha_{l \mid m}^{N}\left(\phi_{N}\right)=Q_{m} \circ P_{m}\left(\mathcal{G}_{l}^{N}\right), \quad m, l=\{1,2, \ldots, L\} .
$$

We study the matrix of "reliabilities",

$$
\begin{equation*}
E_{l \mid m}(\phi)=\varlimsup_{N \rightarrow \infty}-\frac{1}{N} \log \alpha_{l \mid m}^{(N)}\left(\phi_{N}\right), m, l=\{1,2, \ldots, L\} \tag{2}
\end{equation*}
$$

Note that from definitions (1) and (2) it follows that:

$$
\begin{equation*}
E_{m \mid m}=\min _{l \neq m} E_{l \mid m} \tag{3}
\end{equation*}
$$

Definition. The test sequence $\Phi^{*}=\left(\phi_{1}, \phi_{2}, \ldots\right)$ is called LAO if for a given family of positive numbers $E_{1 \mid 1}, E_{2 \mid 2}, \ldots, E_{L-1 \mid L-1}$, the reliability matrix contains in the diagonal these numbers and the remained $L^{2}-L+1$ components take the maximal possible values.

Let $P=\{P(j \mid i)\}$ be a irreducible matrix of transition probabilities of some stationary Markov chain with the same set $\mathcal{X}$ of states, and $Q=\{Q(i), i=\{1,2, \ldots, I\}$ be the corresponding stationary PD.

For a given family of positive numbers $E_{1 \mid 1}, E_{2 \mid 2}, \ldots, E_{L-1 \mid L-1}$, let us define the decision rule $\phi^{*}$ by the following sets

$$
\begin{align*}
\mathcal{R}_{l} & =\left\{Q \circ P: D\left(Q \circ P \| Q \circ P_{l}\right) \leq E_{l \mid l}, \quad D\left(Q \| Q_{l}\right)<\infty\right\}, l=\{1,2, \ldots, L-1\}, \\
\mathcal{R}_{L} & =\left\{Q \circ P: D\left(Q \circ P \| Q \circ P_{l}\right)>E_{l \mid l}, l=\{1,2, \ldots, L-1\},\right.  \tag{4}\\
\mathcal{R}_{l}^{N} & =\mathcal{R}_{l} \cap \mathcal{Q} \circ \mathcal{P}^{N}(\mathcal{X}), \quad l=\{1,2, \ldots, L\} .
\end{align*}
$$

and introduce the functions:

$$
\begin{align*}
& E_{l \mid l}^{*}\left(E_{l \mid l}\right)=E_{l \mid l}, l=\{1,2, \ldots, L-1\}  \tag{5}\\
& E_{l \mid m}^{*}\left(E_{l \mid l}\right)=\inf _{Q \circ P \in \mathcal{R}_{l}} D\left(Q \circ P \| Q \circ P_{m}\right), m=\{1, \ldots, L\}, l \neq m, l=\{1, \ldots, L-1\}, \\
& E_{L \mid m}^{*}\left(E_{1 \mid 1}, \ldots, E_{L-1 \mid L-1}\right)=\inf _{Q \circ P \in \mathcal{R}_{L}} D\left(Q \circ P \| Q \circ P_{m}\right), m=\{1,2, \ldots, L-1\}, \\
& E_{L \mid L}^{*}\left(E_{1 \mid 1}, \ldots, E_{L-1 \mid L-1}\right)=\min _{l=\{1,2, \ldots, L-1\}} E_{l \mid L}^{*} .
\end{align*}
$$

We cite the statement of the general case of large deviation result for types by Natarajan [10].

Theorem 1. Let $\mathcal{X}=\{1,2, \ldots, I\}$ be a finite set of the states of the stationary Markov chain, possessing an irreducible transition matrix $P$ and $\mathcal{A}$ be a nonempty and open subset or convex subset of joint distributions $Q \circ P$ and $Q_{m}$ is stationary distribution for $P_{m}$, then for the type $Q \circ P(\mathbf{x})$ of a vector $\mathbf{x}$ from $Q_{m} \circ P_{m}$ on $\mathcal{X}$ :

$$
\lim _{N \rightarrow \infty}-\frac{1}{N} \log Q_{m} \circ P_{m}^{N}\{\mathbf{x}: Q \circ P(\mathbf{x}) \in \mathcal{A}\}=\inf _{Q \circ P \in \mathcal{A}} D\left(Q \circ P \| Q \circ P_{m}\right) .
$$

Now we formulate the theorem from [7], which we prove by application of Theorem 1.

Theorem 2. Let $\mathcal{X}$ be a fixed finite set, for a family of distinct distributions $P_{1}, \ldots, P_{L}$ the following two statements hold. If the positive finite numbers $E_{1 \mid 1}, \ldots, E_{L-1 \mid L-1}$ satisfy conditions:
$0<E_{1 \mid 1}<\min \left[D\left(Q_{m} \circ P_{m} \| Q_{m} \circ P_{1}\right), m=\{2, \ldots, L\}\right]$,
$0<\quad E_{l \mid l}<\min \left[\begin{array}{l}E_{l \mid m}^{*}\left(E_{m \mid m}\right), m=\{1,2, \ldots, l-1\}, \\ D\left(Q_{m} \circ P_{m} \| Q_{m} \circ P_{l}\right), m=\{l+1, \ldots, L\}\end{array}\right], l=\{2, \ldots, L-1\}$, then:
a) there exists a LAO sequence of tests $\phi^{*}$, the reliability matrix of which $\left\{E_{l \mid m}^{*}\left(\phi^{*}\right)\right\}$ is defined in (5), and all elements of it are positive,
b) even if one of conditions (6) is violated, then the reliability matrix of an arbitrary test necessarily has an element equal to zero, (the corresponding error probability does not tend exponentially to zero).

Proof: First we remark that $D\left(Q \circ P_{l} \| Q \circ P_{m}\right)>0$, for $l \neq m$, because all measures $P_{l}, l=\{1,2, \ldots, L\}$, are distinct. Let us prove the statement a) of the theorem 2 about the existence of the sequence corresponding to a given $E_{1 \mid 1}, \cdots, E_{L-1 \mid L-1}$ satisfying condition (6). Consider the following sequence of tests $\phi^{*}$ given by the sets

$$
\begin{equation*}
\mathcal{B}_{l}^{N}=\bigcup_{Q \circ P \in \mathcal{R}_{l}^{N}} \mathcal{T}_{Q \circ P}^{N}(\mathbf{x}), l=\{1,2, \ldots, L\} . \tag{7}
\end{equation*}
$$

Notice that on account of condition (6) and the continuity of divergence $D$ for $N$ large enough the sets $\mathcal{R}_{l}^{N}, l=\{1,2, \ldots, L\}$ from (4) are not empty. The sets $\mathcal{B}_{l}^{N}, l=$ $\{1,2, \ldots, L\}$, satisfy conditions :

$$
\mathcal{B}_{l}^{N} \bigcap \mathcal{B}_{m}^{N}=\emptyset, l \neq m, \quad \bigcup_{l=1}^{L} \mathcal{B}_{l}^{N}=\mathcal{X}^{N}
$$

Now let us demonstrate that, exponent $E_{l \mid m}\left(\phi^{*}\right)$ for sequence of tests $\phi^{*}$ defined in (7) is equal to $E_{l \mid m}^{*}$. We know from (4) that $\mathcal{R}_{l}, l=\{1,2, \ldots, L-1\}$, are convex subset and $\mathcal{R}_{L}$ is open subset of the decision rule of $\phi^{*}$, therefore $\mathcal{R}_{l}, \quad l=\{1,2, \ldots, L\}$, satisfy in condition of Theorem 1. With relations (4), (5), by Theorem 1 we have

$$
\lim _{N \rightarrow \infty}-\frac{1}{N} \log \alpha_{l \mid m}^{N}\left(\phi^{*}\right)=\lim _{N \rightarrow \infty}-\frac{1}{N} \log Q \circ P_{m}^{N}\left(\mathcal{R}_{l}\right)=\inf _{Q \circ P \in \mathcal{R}_{l}} D\left(Q \circ P \| Q \circ P_{m}\right)
$$

Now using (2) we obtain the following:

$$
E_{l \mid m}\left(\phi^{*}\right)=\inf _{Q \circ P \in \mathcal{R}_{l}} D\left(Q \circ P \| Q \circ P_{m}\right) \quad m, l=\{1,2, \ldots, L\}
$$

Using (6), (4) and (5) it can be realized that all $E_{l \mid m}^{*}$ are strictly positive. The proof of part (a) will be concluded if one demonstrates that the sequence of the tests $\phi^{*}$ is LAO, that is at given finite $E_{1 \mid 1}, \cdots, E_{L-1 \mid L-1}$ for any other sequence of tests $\phi^{* *}$

$$
E_{l \mid m}^{*}\left(\phi^{* *}\right) \leq E_{l \mid m}^{*}\left(\phi^{*}\right), \quad m, l=\{1,2, \ldots, L\}
$$

For this purpose it is sufficient to demonstrate that the sequence of tests will not asymptotically become better if the sets $\mathcal{R}_{l}^{N}$ are not a union of some number of whole types $\mathcal{T}_{Q \circ P}^{N}(\mathbf{x})$, indeed, if a test $\phi^{* *}$ is defined, for instance, by sets $\mathcal{G}_{1}^{N}, \cdots, \mathcal{G}_{L}^{N}$ and, furthermore, $Q \circ P$ is such that

$$
0 \neq Q_{l} \circ P_{l}^{N}\left(\mathcal{G}_{l}^{N} \bigcap \mathcal{T}_{Q \circ P}^{N}(\mathbf{x})\right)=Q_{l} \circ P_{l}^{N}\left(\mathcal{T}_{Q \circ P}^{N}(\mathbf{x})\right)+o(1)
$$

Then, the test $\phi^{* *}$ will be improved, if instead of the set $\mathcal{G}_{l}^{N}$ one takes $\mathcal{G}_{l}^{N} \cup \mathcal{T}_{Q \circ P}^{N}(\mathbf{x})$, as the error probability $\alpha_{l \mid m}^{N}$ can decrease for $m \neq l$. The statement of part (b) of theorem is evident, since the violation of one of the conditions (6) reduces to the equality to zero of at least one of the elements $E_{l \mid m}^{*}$ defined in (5).

## 3. On Statistical Identification of Markov Chain of Distribution Subject to the reliability

Assume that there are $L \geq 2$ hypothetical distributions. The question here is whether or not $r$-th distribution has occurred.

There are two error probabilities for each $r=\{1,2, \ldots, L\}$, the probability $\alpha_{l \neq r \mid m=r}^{(N)}$ to accept $l$ different from $r$, when $r$ is in reality, and the probability $\alpha_{l=r \mid m \neq r}^{(N)}$ that $r$ is accepted, when it is not correct. The probability $\alpha_{l \neq r \mid m=r}^{(N)}$ we already know, it is the probability $\alpha_{r \mid r}^{(N)}$ which is equal to $\sum_{l: l \neq r} \alpha_{l \mid r}^{(N)}$. The reliability $E_{l \neq r \mid m=r}$ coincides with $E_{r \mid r}$, with (3). Our aim is to find the interdependence between $E_{l \neq r \mid m=r}$ and $E_{m \neq r \mid l=r}$. The latter, can have values satisfying in conditions (6), thus, we will have the following conditions:

$$
0<E_{r \mid r}<\min _{l: l \neq r}\left[D\left(Q_{l} \circ P_{l} \| Q_{l} \circ P_{r}\right), r=\{1,2, \ldots, L\} .\right.
$$

We need to use the probabilities of different hypotheses. Let us assume that the hypotheses $H_{l}, l=\{1,2, \ldots, L\}$, have positive probabilities say $P_{\mathbf{r}}(r), r=\{1,2, \ldots, L\}$. We will see that the formulated result in the following theorem, does not depend on values of $P_{\mathbf{r}}(r), r=\{1,2, \ldots, L\}$, if they all are strictly positive. Thus we can make the following calculations for $r=\{1,2, \ldots, L\}$ :

$$
\alpha_{l=r \mid m \neq r}^{(N)}=\frac{P_{\mathbf{r}}^{(N)}(l=r, m \neq r)}{P_{\mathbf{r}}(m \neq r)}=\frac{1}{\sum_{m: m \neq r} P_{\mathbf{r}}(m)} \sum_{m: m \neq r} P_{\mathbf{r}}^{(N)} \alpha_{r \mid m}^{(N)},
$$

and also for $r=\{1,2, \ldots, L\}$, we obtain the following:

$$
\begin{align*}
E_{l=r \mid m \neq r} & =\overline{\lim }_{N \rightarrow \infty}\left(-\frac{1}{N} \log \alpha_{l=r \mid m \neq r}^{(N)}\right)  \tag{8}\\
& =\overline{\lim }_{N \rightarrow \infty} \frac{1}{N}\left(\log \sum_{m: m \neq r} P_{\mathbf{r}}(m)-\log \sum_{m: m \neq r} P_{\mathbf{r}}^{(N)} \alpha_{r \mid m}^{(N)}\right)=\min _{m: m \neq r} E_{r \mid m}^{*} .
\end{align*}
$$

Using (8) by analogy with Theorem 1 and Theorem 2, we conclude (with $\mathcal{R}_{r}$ as in (4) for each $r$ including $r=L$ by the values of $E_{r \mid r}$ from $\left(0, \min D\left(Q_{l} \circ P_{l} \| Q_{l} \circ P_{r}\right)\right)$ ), that

$$
\begin{align*}
& E_{l=r \mid m \neq r}\left(E_{r \mid r}\right)=\min _{m: m \neq r} \inf _{Q \circ P \in \mathcal{R}_{r}} D\left(Q \circ P \| Q \circ P_{m}\right)  \tag{9}\\
& \quad=\min _{m: m \neq r} \quad Q \circ P: D\left(Q \circ P \| Q \circ P_{r}\right) \leq E_{r \mid r} \\
& \quad D\left(Q \circ P \| Q \circ P_{m}\right), r=\{1,2, \ldots, L\} .
\end{align*}
$$

Thus, the obtained result may be formulate in the following theorem.
Theorem 3. For the model with distinct distributions for the given sample $\mathbf{x}$, we can determine its type $Q \circ P$, and when $Q \circ P \in \mathcal{R}_{r}^{(N)}$, we accept the hypotheses $r$. Under the condition that the probabilities of all L hypotheses are positive, the reliability of such test $E_{l=r \mid m \neq r}$ for given $E_{l \neq r \mid m=r}$ is defined in (9).

## 4. On Identification of Two Independent Markov chain of Distributions

In this section, we expand the concept of section 3 for two independent homogenies stationary finite Markov chain. Let $x_{1}$ and $x_{2}$ be independent RV, taking values in the same finite state of Markov chain of set $\mathcal{X}$ with one of $L P D s$, being characteristics of corresponding independent objects, the random vector $\left(X_{1}, X_{2}\right)$ assume values $\left(x^{1}, x^{2}\right) \in \mathcal{X} \times \mathcal{X}$.

Let
$\left(\mathbf{x}_{1}, \mathbf{x}_{\mathbf{2}}\right)=\left(\left(x_{0}^{1}, x_{0}^{2}\right), \ldots,\left(x_{n}^{1}, x_{n}^{2}\right), \ldots,\left(x_{N}^{1}, x_{N}^{2}\right)\right), x^{i} \in \mathcal{X}, i=1,2, n=\{0,2, \ldots, N\}$, be a sequence of results of $N+1$ independent observations of a simple homogeneses stationary Markov chain with finite number $I$ of states. The statistication must define unknown $P D s$ of the objects on the basis of observed data. The selection for each object was done and it was denoted by $\Phi_{N}$. The objects independence test $\Phi_{N}$ may be considered as the pair of tests $\varphi_{N}^{1}$ and $\varphi_{N}^{2}$ for the respective separate objects. We will show the whole compound test sequence by $\Phi$. The test $\varphi_{N}^{i}$ is defined by a partition of the space $\mathcal{X}^{N+1}$ on the $L$ sets and to every trajectory $\mathbf{x}$ the test $\phi_{N}$ puts in one correspondence from $L$ hypotheses. So the space $\mathcal{X}^{N+1}$ will be divided into $L$ parts,

$$
\mathcal{G}_{l, i}^{N}=\left\{\mathbf{x}_{\mathbf{i}}, \phi_{N}\left(\mathbf{x}_{\mathbf{i}}\right)=l\right\}, l=\{1,2, \ldots, L\}, i=1,2 .
$$

We define

$$
\alpha_{l_{1}, l_{2} \mid m_{1}, m_{2}}\left(\Phi_{N}\right)=Q_{m_{1}} \circ P_{m_{1}}\left(\mathcal{G}_{l_{1}, 1}^{N}\right) Q_{m_{2}} \circ P_{m_{2}}\left(\mathcal{G}_{l_{2}, 2}^{N}\right)
$$

as the probability of the erroneous acceptance by the test $\Phi_{N}$ of the hypotheses pair $\left(H_{l_{1}}, H_{l_{2}}\right)$, provided that $\left(H_{m_{1}}, H_{m_{2}}\right)$ is true, where $\left(m_{1}, m_{2}\right) \neq\left(l_{1}, l_{2}\right), m_{i}, l_{i}=$ $\{1,2, \ldots, L\}, i=1,2$. The probability to reject a true pair of hypotheses $\left(H_{m_{1}}, H_{m_{2}}\right)$ by analogy with(1) is the following:

$$
\begin{equation*}
\alpha_{m_{1}, m_{2} \mid m_{1}, m_{2}}^{N}\left(\Phi_{N}\right) \triangleq \sum_{\left(l_{1}, l_{2}\right) \neq\left(m_{1}, m_{2}\right)} \alpha_{l_{1}, l_{2} \mid m_{1}, m_{2}}^{N}\left(\Phi_{N}\right) \tag{10}
\end{equation*}
$$

We also study corresponding limits $E_{l_{1}, l_{2} \mid m_{1}, m_{2}}\left(\Phi_{N}\right)$ of error probability exponents of the sequence of tests $\Phi$, called reliabilities :

$$
\begin{equation*}
E_{l_{1}, l_{2} \mid m_{1}, m_{2}}(\Phi) \triangleq \overline{\lim _{N \rightarrow \infty}}-\frac{1}{N} \log \alpha_{l_{1}, l_{2} \mid m_{1}, m_{2}}\left(\Phi_{N}\right), \quad m_{i}, l_{i}=\overline{1, L}, \quad i=1,2 \tag{11}
\end{equation*}
$$

We denote by $E\left(\varphi^{i}\right)$ the reliability matrices of the sequences of tests $\varphi^{i}, i=1,2$, for each of the objects.

Applying (10) and (11), we obtain the following :

$$
E_{m_{1}, m_{2} \mid m_{1}, m_{2}}(\Phi)=\min _{\left(l_{1}, l_{2}\right) \neq\left(m_{1}, m_{2}\right)} E_{l_{1}, l_{2} \mid m_{1}, m_{2}}(\Phi) .
$$

In this section we use the following lemma.

Lemma. [8], [11] If elements $E_{l \mid m}\left(\varphi^{i}\right), m, l=\{1,2, \ldots, L\}, i=1,2$, are strictly positive, then the following equalities hold for $\Phi=\left(\varphi^{1}, \varphi^{2}\right)$ :

$$
\begin{align*}
& E_{l_{1}, l_{2} \mid m_{1}, m_{2}}(\Phi)=E_{l_{1} \mid m_{1}}\left(\varphi^{1}\right)+E_{l_{2} \mid m_{2}}\left(\varphi^{2}\right), \quad \text { if } \quad m_{1} \neq l_{1}, \quad m_{2} \neq l_{2}  \tag{12a}\\
& E_{l_{1}, l_{2} \mid m_{1}, m_{2}}(\Phi)=E_{l_{i} \mid m_{i}}\left(\varphi^{i}\right), \quad \text { if } \quad m_{3-i}=l_{3-i} \quad m_{i} \neq l_{i}, \quad i=1,2 \tag{12b}
\end{align*}
$$

Consider for given positive elements $E_{m, m \mid m, L}$ and $E_{m, m \mid L, m}, m=\{1,2, \ldots, L-$ $1\}$, the family of regions:

$$
\begin{array}{ll}
\mathcal{R}_{m}^{(1)} \triangleq\left\{Q \circ P: D\left(Q \circ P \| Q \circ P_{m}\right) \leq E_{m, m \mid L, m}\right\}, & m=\{1,2, \ldots, L-1\}, \\
\mathcal{R}_{m}^{(2)} \triangleq\left\{Q \circ P: D\left(Q \circ P \| Q \circ P_{m}\right) \leq E_{m, m \mid m, L}\right\}, & m=\{1,2, \ldots, L-1\}, \\
\mathcal{R}_{L}^{(1)} \triangleq\left\{Q \circ P: D\left(Q \circ P \| Q \circ P_{m}\right)>E_{m, m \mid L, m},\right. & m=\{1,2, \ldots, L-1\}, \\
\mathcal{R}_{L}^{(2)} \triangleq\left\{Q \circ P: D\left(Q \circ P \| Q \circ P_{m}\right)>E_{m, m \mid m, L},\right. & m=\{1,2, \ldots, L-1\} .
\end{array}
$$

What is the identification of the probability distributions for two independent objects? The answer for this question constitutes a reply of the question whether or not the pair of the pair of distributions $\left(r_{1}, r_{2}\right)$ have occurred.

There are two error probabilities for each $\left(r_{1}, r_{2}\right), r_{i}=\{1,2, \ldots, L\}, i=1,2$, the probability $\alpha_{\left(l_{1}, l_{2}\right) \neq\left(r_{1}, r_{2}\right) \mid\left(m_{1}, m_{2}\right)=\left(r_{1}, r_{2}\right)}^{(N)}$ to accept $\left(l_{1}, l_{2}\right)$ different from $\left(r_{1}, r_{2}\right)$, when $\left(r_{1}, r_{2}\right)$ is in reality, and the probability $\alpha_{\left(l_{1}, l_{2}\right)=\left(r_{1}, r_{2}\right) \mid\left(m_{1}, m_{2}\right) \neq\left(r_{1}, r_{2}\right)}^{(N)}$ that $\left(r_{1}, r_{2}\right)$ is accepted, when it is not correct. The probability $\alpha_{\left(l_{1}, l_{2}\right) \neq\left(r_{1}, r_{2}\right) \mid\left(m_{1}, m_{2}\right)=\left(r_{1}, r_{2}\right)}^{(N)}$ is already known, it coincides with the probability $\alpha_{\left(r_{1}, r_{2}\right) \mid\left(r_{1}, r_{2}\right)}^{(N)}$. Our aim is to determine the dependence of $\alpha_{\left(l_{1}, l_{2}\right)=\left(r_{1}, r_{2}\right) \mid\left(m_{1}, m_{2}\right) \neq\left(r_{1}, r_{2}\right)}^{(N)}$ on given $\alpha_{\left(r_{1}, r_{2}\right) \mid\left(r_{1}, r_{2}\right)}^{(N)}$.

We need to use the probabilities of different hypotheses. Let us assume that the hypotheses $H_{l}: l=\{1,2, \ldots, L\}$ have, say, probabilities $P_{\mathbf{r}}(r), r=\{1,2, \ldots, L\}$. The only supposition we shall use is that $P_{\mathbf{r}}(r)>0, r=\{1,2, \ldots, L\}$. We demonstrate, the result formulated in the following theorem does not depend on values of $P_{\mathbf{r}}(r), r=$ $\{1,2, \ldots, L\}$, if they all are strictly positive. Thus, the following reasoning can be made for each $r_{i}=\{1,2, \ldots, L\}, i=1,2$ :

$$
\begin{aligned}
& \alpha_{\left(l_{1}, l_{2}\right)=\left(r_{1}, r_{2}\right) \mid\left(m_{1}, m_{2}\right) \neq\left(r_{1}, r_{2}\right)}^{(N)}=\frac{P_{\mathbf{r}}^{(N)}\left(\left(l_{1}, l_{2}\right)=\left(r_{1}, r_{2}\right),\left(m_{1}, m_{2}\right) \neq\left(r_{1}, r_{2}\right)\right)}{\left.P_{\mathbf{r}}\left(m_{1}, m_{2}\right) \neq\left(r_{1}, r_{2}\right)\right)}, \\
& \alpha_{\left(l_{1}, l_{2}\right)=\left(r_{1}, r_{2}\right) \mid\left(m_{1}, m_{2}\right) \neq\left(r_{1}, r_{2}\right)}^{(N)} \underset{m:\left(m_{1}, m_{2}\right) \neq\left(r_{1}, r_{2}\right)}{\sum_{\mathbf{r}}\left(m_{1}, m_{2}\right)} \sum_{m:\left(m_{1}, m_{2}\right) \neq\left(r_{1}, r_{2}\right)} \alpha_{\left(m_{1}, m_{2}\right) \mid\left(r_{1}, r_{2}\right)} P_{\mathbf{r}}^{(N)}\left(m_{1}, m_{2}\right) .
\end{aligned}
$$

Finally, for $r=\overline{1, L}$, we obtain the following:

$$
\begin{equation*}
E_{\left(l_{1}, l_{2}\right)=\left(r_{1}, r_{2}\right) \mid\left(m_{1}, m_{2}\right) \neq\left(r_{1}, r_{2}\right)}=\min _{\left(m_{1}, m_{2}\right):\left(m_{1}, m_{2}\right) \neq\left(r_{1}, r_{2}\right)} E_{\left(r_{1}, r_{2}\right) \mid\left(m_{1}, m_{2}\right)}^{*} . \tag{13}
\end{equation*}
$$

For every LAO test $\Phi^{*}$ from (11), (12) and (13) we obtain the following:

$$
E_{\left(l_{1}, l_{2}\right)=\left(r_{1}, r_{2}\right) \mid\left(m_{1}, m_{2}\right) \neq\left(r_{1}, r_{2}\right)}=\min _{m_{1} \neq r_{1}, m_{2} \neq r_{2}}\left(E_{r_{1} \mid m_{1}}^{1}, E_{r_{2} \mid m_{2}}^{2}\right)
$$

where, $E_{r_{1} \mid m_{1}}^{1}, E_{r_{2} \mid m_{2}}^{2}$ are determined by (5) for, correspondingly, the first and the second objects. For every LAO test $\Phi^{*}$ from (11) and (12) we deduce that

$$
\begin{equation*}
E_{\left(r_{1}, r_{2}\right) \mid\left(r_{1}, r_{2}\right)}=\min _{m_{1} \neq r_{1}, m_{2} \neq r_{2}}\left(E_{r_{1} \mid m_{1}}^{1}, E_{r_{2} \mid m_{2}}^{2}\right)=\min \left(E_{r_{1} \mid r_{1}}^{1}, E_{r_{2} \mid r_{2}}^{2}\right) \tag{14}
\end{equation*}
$$

and each of $E_{r_{1} \mid r_{1}}^{1}, E_{r_{2} \mid r_{2}}^{2}$ satisfy the following conditions (see theorem 2, condition (6)).

$$
\begin{align*}
& 0<E_{r_{1} \mid r_{1}}^{1}<\min \left[\min _{l=\left\{1, \ldots, r_{1}-1\right\}} E_{l \mid m}^{*}\left(E_{l \mid l}^{1}\right), \min _{l=\left\{r_{1}+1, \ldots, L\right\}} D\left(Q_{l} \circ P_{l} \| Q_{l} \circ P_{r_{1}}\right)\right],  \tag{15a}\\
& 0<E_{r_{2} \mid r_{2}}^{2}<\min \left[\min _{l=\left\{1,2, \ldots, r_{2}-1\right\}} E_{l \mid m}^{*}\left(E_{l \mid l}^{2}\right), \min _{l=\left\{r_{2}+1, \ldots, L\right\}} D\left(Q_{l} \circ P_{l} \| Q_{l} \circ P_{r_{2}}\right)\right], \tag{15b}
\end{align*}
$$

From (5), we see that the elements $E_{l \mid m}^{*}\left(E_{l \mid l}^{1}\right), r_{1}=\left\{1,2, \ldots, r_{1}-1\right\}$ and $E_{l \mid m}^{*}\left(E_{l \mid l}^{2}\right), r_{2}=$ $\left\{1,2, \ldots, r_{2}-1\right\}$ are determined by only $E_{l \mid l}^{1}$ and $E_{l \mid l}^{2}$. However, we are considering only elements $E_{r_{1} \mid r_{1}}^{1}$ and $E_{r_{2} \mid r_{2}}^{2}$. Using theorem 1 and (15) we have

$$
\begin{align*}
0 & <E_{r_{1} \mid r_{1}}^{1} \\
& <\min \left[\min _{l=\left\{1,2, \ldots, r_{1}-1\right\}} D\left(Q_{l} \circ P_{l} \| Q_{l} \circ P_{r_{1}}\right), \min _{l=\left\{r_{1}+1, \ldots, L\right\}} D\left(Q_{l} \circ P_{l} \| Q_{l} \circ P_{r_{1}}\right)\right],  \tag{16a}\\
0 & <E_{r_{2} \mid r_{2}}^{2} \\
& <\min \left[\min _{l=\left\{1,2, \ldots, r_{2}-1\right\}} D\left(Q_{l} \circ P_{l} \| Q_{l} \circ P_{r_{2}}\right), \min _{l=\left\{r_{2}+1, \ldots, L\right\}} D\left(Q_{l} \circ P_{l} \| Q_{l} \circ P_{r_{2}}\right)\right] . \tag{16b}
\end{align*}
$$

Let us denote $r=\max \left(r_{1}, r_{2}\right)$ and $k=\min \left(r_{1}, r_{2}\right)$. From (14) we have that, when $E_{\left(r_{1}, r_{2}\right) \mid\left(r_{1}, r_{2}\right)}=E_{r_{1} \mid r_{1}}^{1}$, then $E_{r_{1} \mid r_{1}}^{1} \leq E_{r_{2} \mid r_{2}}^{2}$ and when $E_{\left(r_{1}, r_{2}\right) \mid\left(r_{1}, r_{2}\right)}=E_{r_{2} \mid r_{2}}^{2}$, then $E_{r_{1} \mid r_{1}}^{1} \geq E_{r_{2} \mid r_{2}}^{2}$. Thus, it can be implied that given strictly positive elements $E_{\left(r_{1}, r_{2}\right) \mid\left(r_{1}, r_{2}\right)}$ must meet both inequalities (16) and the combination of these restrictions gives

$$
\begin{align*}
0< & E_{\left(r_{1}, r_{2}\right) \mid\left(r_{1}, r_{2}\right)} \\
& <\min \left[\min _{l=\{1,2, \ldots,, r-1\}} D\left(Q_{l} \circ P_{l} \| Q_{l} \circ P_{r}\right), \min _{l=\{r+1, \ldots, L\}} D\left(Q_{l} \circ P_{l} \| Q_{l} \circ P_{k}\right)\right] . \tag{17}
\end{align*}
$$

Using (15) we can determine reliability $E_{\left(l_{1}, l_{2}\right)=\left(r_{1}, r_{2}\right) \mid\left(m_{1}, m_{2}\right) \neq\left(r_{1}, r_{2}\right)}$ in function of $E_{\left(r_{1}, r_{2}\right) \mid\left(r_{1}, r_{2}\right)}$ as follows:

$$
\begin{align*}
& E_{\left(l_{1}, l_{2}\right)=\left(r_{1}, r_{2}\right) \mid\left(m_{1}, m_{2}\right) \neq\left(r_{1}, r_{2}\right)}\left(E_{\left(r_{1}, r_{2}\right) \mid\left(r_{1}, r_{2}\right)}\right) \\
&=\min _{m_{1} \neq r_{1}, m_{2} \neq r_{2}}\left(E_{r_{1} \mid r_{1}}\left(E_{\left(r_{1}, r_{2}\right) \mid\left(r_{1}, r_{2}\right)}\right), E_{r_{2} \mid r_{2}}\left(E_{\left(r_{1}, r_{2}\right) \mid\left(r_{1}, r_{2}\right)}\right)\right. \tag{18}
\end{align*}
$$

where, $\left(E_{r_{1} \mid r_{1}}\left(E_{\left(r_{1}, r_{2}\right) \mid\left(r_{1}, r_{2}\right)}\right)\right.$ and $E_{r_{2} \mid r_{2}}\left(E_{\left(r_{1}, r_{2}\right) \mid\left(r_{1}, r_{2}\right)}\right)$ are determined by (5). The obtained results can be summarized in the following theorem:

Theorem 4. If the distributions $H_{m}, m=\{1,2, \ldots, L\}$, are different and the given strictly positive numbers $E_{\left(r_{1}, r_{2}\right) \mid\left(r_{1}, r_{2}\right)}$ satisfies condition (17), then the reliability $E_{\left(l_{1}, l_{2}\right)=\left(r_{1}, r_{2}\right) \mid\left(m_{1}, m_{2}\right) \neq\left(r_{1}, r_{2}\right)}$ is defined in (18).

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# On intuitionistic fuzzy ideals in $\Gamma$ - rings 

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#### Abstract

In this paper, we study some properties of intuitionistic fuzzy ideals of a $\Gamma$ - ring and prove some results on these.


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Key Words and Phrases: $\Gamma$ - ring, fuzzy set, intuitionistic fuzzy set, intuitionistic fuzzy ideal

## 1 Introduction

The notion of a fuzzy set was introduced by L.A.Zadeh [12], and since then this concept has been applied to various algebraic structures. The idea of "Intuitionistic Fuzzy Set" was first introduced by K.T.Atanassov [1] as a generalization of the notion of fuzzy set. N.Nobusawa [10] introduced the notion of a $\Gamma$ - ring, as more general than a ring. W.E.Barnes[2] weakened slightly the conditions in the definition of the $\Gamma$ - rings in the sense of Nobusawa. W.E.Barnes [2], S.Kyuno [7,8] and J.Luh [9] studied the structure of $\Gamma$ - rings and obtained various generalizations analogous to corresponding parts in ring theory. Y.B.Jun and C.Y.Lee [5] introduced the concept of fuzzy ideals in the theory of $\Gamma$ - rings. In this paper, we study the notion of intuitionistic fuzzy ideals in $\Gamma$ - rings and prove some of its properties.

## 2 Preliminaries

In this section the definition of $\Gamma$-ring in the sense of Nobusawa and Barnes is discussed with examples. Also we include some elementary concepts that are necessary for this paper. Professor N.Nobusawa introduced the concept of $\Gamma$-ring and Professor W.Barnes generalized this concept.

Definition 2.1[2]. If $M=\{x, y, z, \cdots\}$ and $\Gamma=\{\alpha, \beta, \gamma \cdots\}$ be two additive abelian groups and for all $x, y, z \in M$ and $\alpha, \beta \in \Gamma$, the following conditions are satisfied:

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1. $x \alpha y \in M$,
2. $(x+y) \alpha z=x \alpha z+y \alpha z, x(\alpha+\beta) y=x \alpha y+x \beta y, x \alpha(y+z)=x \alpha y+x \alpha z$,
3. $(x \alpha y) \beta z=x \alpha(y \beta z)$, then $M$ is called a $\Gamma$ - ring.

If these conditions are strengthened to
$\left(1^{\prime}\right) x \alpha y \in M, \alpha x \beta \in \Gamma$,
$\left(2^{\prime}\right)(x+y) \alpha z=x \alpha z+y \alpha z, x(\alpha+\beta) y=x \alpha y+x \beta y, x \alpha(y+z)=x \alpha y+x \alpha z$,
$\left(3^{\prime}\right)(x \alpha y) \beta z=x(\alpha y \beta) z=x \alpha(y \beta z)$,
(4') $x \alpha y=0$ for all $x, y \in M$ implies $\alpha=0$,
we then have a $\Gamma$ - ring in the sense of Nobusawa [10]. As indicated in [10], an example of a $\Gamma$ - ring is obtained by letting $X$ and $Y$ be abelian groups, $M=\operatorname{Hom}(X, Y)$, $\Gamma=\operatorname{Hom}(Y, X)$ and $x \alpha y$ be the usual composite map. (While Nobusawa does not explicitly require that $M$ and $\Gamma$ be abelian groups, it appears clear that this is intended.) We may note that it follows from (1)-(3) that $0 \alpha y=x 0 y=x \alpha 0=0$ for all $x, y \in M$ and all $\alpha \in \Gamma$.
Example 2.2. If $G$ and $G^{\prime}$ are two additive abelian groups, $M=\operatorname{Hom}\left(G, G^{\prime}\right), \Gamma=$ $\operatorname{Hom}\left(\mathrm{G}^{\prime}, \mathrm{G}\right)$ then M is a $\Gamma$-ring with respect to point wise addition and composition of mappings.
Example 2.3. If U and V be vector spaces over the same field $\mathrm{F}, \mathrm{M}=\operatorname{Hom}(\mathrm{U}, \mathrm{V})$, $\Gamma=\operatorname{Hom}(\mathrm{V}, \mathrm{U})$. Then M is a $\Gamma$-ring with respect to point wise addition and composition of mappings.
Definition 2.4[2]. A subset $A$ of a $\Gamma$ - ring $M$ is a left (resp. right) ideal of $M$ if $A$ is an additive subgroup of $M$ such that $M \Gamma A \subseteq A$ (resp. $A \Gamma M \subseteq A$ ), where $M \Gamma A=\{x \alpha y \mid x \in M, \alpha \in \Gamma, y \in A\}$ and $A \Gamma M=\{y \alpha x \mid y \in A, \alpha \in \Gamma, x \in M\}$. If $A$ is both a left and a right ideal, than $A$ is a two sided ideal or simply an ideal of $M$.

Definition 2.5[11]. A fuzzy set $A$ in $M$ is a function $A: M \rightarrow[0,1]$.
Definition 2.6[11]. Let $\mu$ be a fuzzy set in a $\Gamma$ - ring $M$. For any $t \in[0,1]$, the set $U(\mu, t)=\{x \in M \mid \mu(x) \geq t\}$ is called a level set of $\mu$.
Definition 2.7[11]. A fuzzy set $\mu$ in a $\Gamma$ - ring $M$ is called a fuzzy left (resp. right) ideal of $M$, if it satisfies:
(i) $\mu(x-y) \geq \mu(x) \wedge \mu(y)$,
(ii) $\mu(x \alpha y) \geq \mu(y)($ resp. $\mu(x \alpha y) \geq \mu(x))$,
for all $x, y \in M$ and $\alpha \in \Gamma$. If $\mu$ is both a fuzzy left and right ideal of $M$, then $\mu$ is called a fuzzy ideal of $M$.

Definition 2.8[1]. Let $X$ be a nonempty fixed set. An intuitionistic fuzzy set $A$ in $X$ is an object having the form $A=\left\{\left\langle x, \mu_{A}(x), \nu_{A}(x)\right\rangle \mid x \in X\right\}$, where the functions $\mu_{A}: X \rightarrow[0,1]$ and $\nu_{A}: X \rightarrow[0,1]$ denote the degree of membership and the degree of non-membership of each element $x \in X$ to the set $A$, respectively, and $0 \leq \mu_{A}(x)+\nu_{A}(x) \leq 1$, for every $x \in X$.

Notation. For the sake of simplicity, we shall denote the intuitionistic fuzzy set (IFS in short) $A=\left\{\left\langle x, \mu_{A}(x), \nu_{A}(x)\right\rangle \mid x \in X\right\}$ by $A=\left\langle\mu_{A}, \nu_{A}\right\rangle$.

Definition 2.9[1]. Let $X$ be a non empty set and let $A=\left\langle\mu_{A}, \nu_{A}\right\rangle$ and $B=\left\langle\mu_{B}, \nu_{B}\right\rangle$ be IFSs in $X$. Then

1. $A \subset B$ iff $\mu_{A} \leq \mu_{B}$ and $\nu_{A} \geq \nu_{B}$.
2. $A=B$ iff $A \subset B$ and $B \subset A$.
3. $A^{c}=\left\langle\nu_{A}, \mu_{A}\right\rangle$.
4. $A \cap B=\left(\mu_{A} \wedge \mu_{B}, \nu_{A} \vee \nu_{B}\right)$.
5. $A \cup B=\left(\mu_{A} \vee \mu_{B}, \nu_{A} \wedge \nu_{B}\right)$.
6.$A=\left(\mu_{A}, 1-\mu_{A}\right), \diamond A=\left(1-\nu_{A}, \nu_{A}\right)$.

Definition 2.10[5]. Let $A$ be an IFS in a $\Gamma$ - ring $M$. For each pair $\langle t, s\rangle \in[0,1]$ with $t+s \leq 1$, the set $A_{\langle t, s\rangle}=\left\{x \in X \mid \mu_{A}(x) \geq t\right.$ and $\left.\nu_{A}(x) \leq s\right\}$ is called a $\langle t, s\rangle$ level subset of $A$.

Definition 2.11.[6]. Let $A=\left\langle\mu_{A}, \nu_{A}\right\rangle$ be an intuitionistic fuzzy set in a $\Gamma$ ring M and let $t \in[0,1]$. Then the sets $U\left(\mu_{A} ; t\right)=\left\{x \in M: \mu_{A}(x) \geq t\right\}$ and $L\left(\nu_{A} ; t\right)=\left\{x \in M: \nu_{A}(x) \leq t\right\}$ are called upper level set and lower level set of $A$ respectively.

## 3 Intuitionistic fuzzy ideals

In what follows, let $M$ denote a $\Gamma$ - ring unless otherwise specified. In this section, an example of an intuitionistic fuzzy ideal is given.
Definition 3.1. An IFS $A=\left\langle\mu_{A}, \nu_{A}\right\rangle$ in $M$ is called an intuitionistic fuzzy left (resp. right) ideal of a $\Gamma$ - ring $M$ if
(i) $\mu_{A}(x-y) \geq\left\{\mu_{A}(x) \wedge \mu_{A}(y)\right\}$ and $\mu_{A}(x \alpha y) \geq \mu_{A}(y)\left(\right.$ resp. $\left.\mu_{A}(x \alpha y) \geq \mu_{A}(x)\right)$,
(ii) $\nu_{A}(x-y) \leq\left\{\nu_{A}(x) \vee \nu_{A}(y)\right\}$ and $\nu_{A}(x \alpha y) \leq \nu_{A}(y)\left(\operatorname{resp} . \nu_{A}(x \alpha y) \leq \nu_{A}(x)\right)$,
for all $x, y \in M$ and $\alpha \in \Gamma$.
Example 3.2. [Intuitionistic fuzzy ideal of a $\Gamma$-ring]
Let $R$ be the set of all integers. Then $R$ is a ring.
Take $\mathrm{M}=\Gamma=\mathrm{R}$.
Let $a, b \in M, \alpha \in \Gamma$. Suppose $a \alpha b$ is the product of $a, \alpha, b \in M$. Then M is a $\Gamma$-ring. Define an IFS A $=\left\langle\mu_{A}, \nu_{A}\right\rangle$ in M as follows.
$\mu_{A}(0)=1$ and $\mu_{A}( \pm 1)=\mu_{A}( \pm 2)=\ldots . .=\mathrm{t}$ and
$\nu_{A}(0)=0$ and $\nu_{A}( \pm 1)=\nu_{A}( \pm 2)=\ldots .=\mathrm{s}$, where $t \in[0,1], s \in[0,1]$ and $t+s \leq 1$. By routine calculations, clearly A is an intuitionistic fuzzy ideal of a $\Gamma$-ring M.

Theorem 3.3. If $A$ is an ideal of a $\Gamma$ - ring $M$, then the IFS $\hat{A}=\left\langle\chi_{A}, \overline{\chi_{A}}\right\rangle$ is an intuitionistic fuzzy ideal of $M$.
Proof. Let $x, y \in M$.
If $x, y \in A$ and $\alpha \in \Gamma$, then $x-y \in A$ and $x \alpha y \in A$, since $A$ is an ideal of $M$.
Hence $\chi_{A}(x-y)=1 \geq\left\{\chi_{A}(x) \wedge \chi_{A}(y)\right\}$ and $\chi_{A}(x \alpha y)=1 \geq \chi_{A}(y)$.
Also, we have
$0=1-\chi_{A}(x-y)=\overline{\chi_{A}}(x-y) \leq\left\{\overline{\chi_{A}}(x) \vee \overline{\chi_{A}}(y)\right\}$ and
$0=1-\chi_{A}(x \alpha y)=\overline{\chi_{A}}(x \alpha y) \leq \overline{\chi_{A}}(y)$.
If $x \notin A$ or $y \notin A$, then $\chi_{A}(x)=0$ or $\chi_{A}(y)=0$. Thus wwe have
$\chi_{A}(x-y) \geq\left\{\chi_{A}(x) \wedge \chi_{A}(y)\right\}$ and
$c h i_{A}(x \alpha y) \geq \chi_{A}(y)$ for all $\alpha \in \Gamma$.
Also $\overline{\chi_{A}}(x-y) \leq\left\{\overline{\chi_{A}}(x) \vee\right.$
overline $\left.\chi_{A}(y)\right\}=\left(1-\chi_{A}(x)\right) \vee\left(1-\chi_{A}(y)\right)=1$
and $\overline{\chi_{A}}(x \alpha y)=1-\chi_{A}(x \alpha y) \leq 1-\chi_{A}(y)=\overline{\chi_{A}}(y)$.
This completes the proof.
Definition 3.4[3]. An intuitionistic fuzzy left (resp. right) ideal $A=\left\langle\mu_{A}, \nu_{A}\right\rangle$ of a $\Gamma$ - ring $M$ is said to be normal if $\mu_{A}(0)=1$ and $\nu_{A}(0)=0$.
Theorem 3.5. Let $A=\left\langle\mu_{A}, \nu_{A}\right\rangle$ be an intuitionistic fuzzy left (resp. right) ideal of a $\Gamma$-ring $M$ and let $\mu_{A}^{+}(x)=\mu_{A}(x)+1-\mu_{A}(0), \nu_{A}^{+}(x)=\nu_{A}(x)-\nu_{A}(0)$. If $\mu_{A}^{+}(x)+\nu_{A}^{+}(x) \leq 1$ for all $x \in M$, then $A^{+}=\left\langle\mu_{A}^{+}, \nu_{A}^{+}\right\rangle$is a normal intuitionistic fuzzy left (resp. right) ideal of $M$.
Proof. We first observe that $\mu_{A}^{+}(0)=1, \nu_{A}^{+}(0)=0$ and $\mu_{A}^{+}(x), \nu_{A}^{+}(x) \in[0,1]$ for every $x \in M$. So, $A^{+}=\left\langle\mu_{A}^{+}, \nu_{A}^{+}\right\rangle$is a normal intuitionistic fuzzy set.
To prove that it is an intuitionistic fuzzy left (resp. right) ideal, let $x, y \in M$ and $\alpha \in \Gamma$.

Then $\mu_{A}^{+}(x-y)=\mu_{A}(x-y)+1-\mu_{A}(0) \geq\left\{\mu_{A}(x) \wedge \mu_{A}(y)\right\}+1-\mu_{A}(0)$

$$
=\left\{\mu_{A}(x)+1-\mu_{A}(0)\right\} \wedge\left\{\mu_{A}(y)+1-\mu_{A}(0)\right\}=\mu_{A}^{+}(x) \wedge \mu_{A}^{+}(y)
$$

$$
\nu_{A}^{+}(x-y)=\nu_{A}(x-y)-\nu_{A}(0) \leq\left\{\nu_{A}(x) \vee \nu_{A}(y)\right\}-\nu_{A}(0)
$$

$$
=\left\{\nu_{A}(x)-\nu_{A}(0)\right\} \vee\left\{\nu_{A}(y)-\nu_{A}(0)\right\}=\nu_{A}^{+}(x) \vee \nu_{A}^{+}(y)
$$

and $\mu_{A}^{+}(x \alpha y)=\mu_{A}(x \alpha y)+1-\mu_{A}(0) \geq \mu_{A}(y)+1-\mu_{A}(0)=\mu_{A}^{+}(y)$,

$$
\nu_{A}^{+}(x \alpha y)=\nu_{A}(x \alpha y)-\nu_{A}(0) \leq \nu_{A}(y)-\nu_{A}(0)=\nu_{A}^{+}(y)
$$

This shows that $A^{+}$is an intuitionistic fuzzy left(resp. right) ideal of $M$. So, $A^{+}=\left\langle\mu_{A}^{+}, \nu_{A}^{+}\right\rangle$is a normal intuitionistic fuzzy left(resp. right)ideal of $M$.
Definition 3.6 [2]. Let $I$ be an ideal of a $\Gamma$ - ring $M$. If for each $a+I, b+I$ in the factor group $M / I$ and each $\alpha \in \Gamma$, we define $(a+I) \alpha(b+I)=a \alpha b+I$, then $M / I$ is a $\Gamma$ - ring which we shall call the $\Gamma$ - residue class ring of $M$ with respect to $I$.
Theorem 3.7. Let $I$ be an ideal of a $\Gamma$ - ring $M$. If $A$ is an intuitionistic fuzzy left (resp. right) ideal of $M$, then the intuitionistic fuzzy set $\tilde{A}$ of $M / I$ defined by

$$
\mu_{\tilde{A}}(a+I)=\bigvee_{x \in I} \mu_{A}(a+x) \text { and } \nu_{\tilde{A}}(a+I)=\bigwedge_{x \in I} \nu_{A}(a+x)
$$

is an intuitionistic fuzzy left (resp. right) ideal of the $\Gamma$ - residue class ring $M / I$ of M with respect to $I$.
Proof. Let $a, b \in M$ be such that $a+I=b+I$.
Then $b=a+y$ for some $y \in I$ and so

$$
\begin{aligned}
& \mu_{\tilde{A}}(b+I)=\bigvee_{x \in I} \mu_{A}(b+x)=\bigvee_{x \in I} \mu_{A}(a+y+x)=\bigvee_{x+y=z \in I} \mu_{A}(a+z)=\mu_{\tilde{A}}(a+I), \\
& \nu_{\tilde{A}}(b+I)=\bigwedge_{x \in I} \nu_{A}(b+x)=\bigwedge_{x \in I} \nu_{A}(a+y+x)=\bigwedge_{x+y=z \in I} \nu_{A}(a+z)=\nu_{\tilde{A}}(a+I)
\end{aligned}
$$

Hence $\tilde{A}$ is well defined.
For any $x+I, y+I \in M / I$ and $\alpha \in \Gamma$, we have

$$
\begin{aligned}
\mu_{\tilde{A}}((x+I)-(y+I))=\mu_{\tilde{A}}((x-y)+I)=\bigvee_{z \in I} \mu_{A} & ((x-y)+z) \\
& =\bigvee_{z=u-v \in I} \mu_{A}((x-y)+(u-v)) \\
& =\bigvee_{u, v \in I} \mu_{A}((x+u)-(y+v)) \\
& \geq \bigvee_{u, v \in I}\left(\mu_{A}(x+u) \wedge \mu_{A}(y+v)\right) \\
& =\left(\bigvee_{u \in I} \mu_{A}(x+u)\right) \wedge\left(\bigvee_{v \in I} \mu_{A}(y+v)\right) \\
& =\mu_{\tilde{A}}(x+I) \wedge \mu_{\tilde{A}}(y+I)
\end{aligned}
$$

$$
\begin{aligned}
& =\bigwedge_{z=u-v \in I} \nu_{A}((x-y)+(u-v)) \\
& =\bigwedge_{u, v \in I} \nu_{A}((x+u)-(y+v)) \\
& \leq \bigwedge_{u, v \in I}\left(\nu_{A}(x+u) \vee \nu_{A}(y+v)\right) \\
& =\left(\bigwedge_{u \in I} \nu_{A}(x+u)\right) \vee\left(\bigwedge_{v \in I} \nu_{A}(y+v)\right) \\
& =\nu_{\tilde{A}}(x+I) \vee \nu_{\tilde{A}}(y+I),
\end{aligned}
$$

$$
\mu_{\tilde{A}}((x+I) \alpha(y+I))=\mu_{\tilde{A}}((x \alpha y)+I)=\bigvee_{z \in I} \mu_{A}((x \alpha y)+z)
$$

$$
\geq \bigvee_{z \in I} \mu_{A}(x \alpha y+x \alpha z) \text { because } x \alpha z \in I
$$

$$
=\bigvee_{z \in I} \mu_{A}(x \alpha(y+z))
$$

$$
\geq \bigvee_{z \in I} \mu_{A}(y+z)
$$

$$
=\mu_{\tilde{A}}(y+I),
$$

$$
\nu_{\tilde{A}}((x+I) \alpha(y+I))=\nu_{\tilde{A}}((x \alpha y)+I)=\bigwedge_{z \in I} \nu_{A}((x \alpha y)+z)
$$

$$
\leq \bigwedge_{z \in I} \nu_{A}(x \alpha y+x \alpha z) \text { because } x \alpha z \in I
$$

$$
=\bigwedge_{z \in I} \nu_{A}(x \alpha(y+z))
$$

$$
\leq \wedge_{z \in I} \nu_{A}(y+z)
$$

$$
=\nu_{\tilde{A}}(y+I) .
$$

Similarly,

$$
\mu_{\tilde{A}}((x+I) \alpha(y+I)) \geq \mu_{\tilde{A}}(x+I) \text { and } \nu_{\tilde{A}}((x+I) \alpha(y+I)) \leq \nu_{\tilde{A}}(x+I) .
$$

Hence $\tilde{A}$ is an intuitionistic fuzzy left (resp. right) ideal of $M / I$.
Theorem 3.8. If the IFS $A=\left\langle\mu_{A}, \nu_{A}\right\rangle$ is an intuitionistic fuzzy left (resp. right) ideal of $M$, then the set $M_{A}=\left\{x \in M \mid \mu_{A}(x)=\mu_{A}(0)\right.$ and $\left.\nu_{A}(x)=\nu_{A}(0)\right\}$ is an ideal of $M$.

Proof. Let $x, y \in M_{A}$.

Then $\mu_{A}(x)=\mu_{A}(y)=\mu_{A}(0)$ and $\nu_{A}(x)=\nu_{A}(y)=\nu_{A}(0)$.
Since $A$ is an intuitionistic fuzzy ideal of $M$, it follows that
$\mu_{A}(x-y) \geq\left\{\mu_{A}(x) \wedge \mu_{A}(y)\right\}=\left\{\mu_{A}(0) \wedge \mu_{A}(0)\right\}=\mu_{A}(0)$,
$\nu_{A}(x-y) \leq\left\{\nu_{A}(x) \vee \nu_{A}(y)\right\}=\left\{\nu_{A}(0) \vee \nu_{A}(0)\right\}=\nu_{A}(0)$.
Hence $\mu_{A}(x-y)=\mu_{A}(0)$ and $\nu_{A}(x-y)=\nu_{A}(0)$. So $x-y \in M_{A}$.
Let $x \in M, \alpha \in \Gamma$ and $y \in M_{A}$.
Therefore $\mu_{A}(x \alpha y) \geq \mu_{A}(y)=\mu_{A}(0)$ (resp. $\left.\mu_{A}(x \alpha y) \geq \mu_{A}(x)=\mu_{A}(0)\right)$ and
$\nu_{A}(x \alpha y) \leq \nu_{A}(y)=\nu_{A}(0)\left(\right.$ resp. $\left.\nu_{A}(x \alpha y) \leq \nu_{A}(x)=\nu_{A}(0)\right)$.
Hence $\mu_{A}(x \alpha y)=\mu_{A}(0)$ and $\nu_{A}(x \alpha y)=\nu_{A}(0)$.
So x $\alpha \mathrm{y} \in M_{A}$. Hence $M_{A}$ is an ideal of $M$.
Theorem 3.9. Let $A$ be an intuitionistic fuzzy left (resp. right) ideal of a $\Gamma$ - ring $M$. For each pair $\langle t, s\rangle \in[0,1]$, the level set $A_{\langle t, s\rangle}$ is an ideal of $M$.

Proof. Let $x, y \in A_{\langle t, s\rangle}$.
Then $\mu_{A}(x) \geq t, \mu_{A}(y) \geq t$ and $\nu_{A}(x) \leq s, \nu_{A}(y) \leq s$.
Since $A$ is an intuitionistic fuzzy left (resp. right) ideal, we have
$\mu_{A}(x-y) \geq\left\{\mu_{A}(x) \wedge \mu_{A}(y)\right\} \geq t$ and $\nu_{A}(x-y) \leq\left\{\nu_{A}(x) \vee \nu_{A}(y)\right\} \leq s$.
So $x-y \in A_{\langle t, s\rangle}$.
Let $x \in M, y \in A_{\langle t, s\rangle}$ and $\alpha \in \Gamma$.
Then $\mu_{A}(x \alpha y) \geq \mu_{A}(y) \geq t$ and $\nu_{A}(x \alpha y) \leq \nu_{A}(y) \leq s$. So $x \alpha y \in A_{\langle t, s\rangle}$.
Hence $A_{\langle t, s\rangle}$ is an ideal of $M$.
Definition 3.10. Let $A$ and $B$ be two intuitionistic fuzzy subsets of a $\Gamma$ - ring $M$ and $\alpha \in \Gamma$. Then the product $A \Gamma B$ is defined by

$$
\begin{aligned}
& \mu_{A \Gamma B}(x)=\bigvee_{x=y \alpha z}\left(\mu_{A}(y) \wedge \mu_{A}(z)\right) \text { if } x=y \alpha z \\
& \nu_{A \Gamma B}(x)=\bigwedge_{x=y \alpha z}\left(\nu_{A}(y) \vee \nu_{A}(z)\right) \text { if } x=y \alpha z
\end{aligned}
$$

Otherwise, we define $\mu_{A \Gamma B}(x)=0$ and $\nu_{A \Gamma B}(x)=1$.
Definition 3.11[4]. Let $A=\left\langle\mu_{A}, \nu_{A}\right\rangle$ and $B=\left\langle\mu_{B}, \nu_{B}\right\rangle$ be two IFSs in a $\Gamma$ - ring $M$. Then the composition of $A$ and $B$ is defined to be the intuitionistic fuzzy set $A \circ B=\left\langle\mu_{A \circ B}, \nu_{A \circ B}\right\rangle$ in $M$ given by

$$
\begin{aligned}
& \mu_{A \circ B}(x)=\bigvee\left\{\begin{array}{l}
\bigwedge_{1 \leq i \leq k} \mu_{A}\left(a_{i}\right) \wedge \mu_{B}\left(b_{i}\right): x=\sum_{1}^{k} a_{i} \alpha b_{i}, a_{i}, b_{i} \in M, \alpha \in \Gamma, k \in N \\
\nu_{A \circ B}(x)=\bigwedge \nu_{A}\left(a_{i}\right) \vee \nu_{B}\left(b_{i}\right): x=\sum_{1}^{k} a_{i} \alpha b_{i}, a_{i}, b_{i} \in M, \alpha \in \Gamma, k \in N
\end{array}\right.
\end{aligned}
$$

if we can express $x=\sum_{i=1}^{k} a_{i} \alpha b_{i}$ for some $a_{i}, b_{i} \in M$, where each $a_{i} \alpha b_{i} \neq 0$ and $k \in N$.

Otherwise, we define $A \circ B=0$, i.e., $\mu_{A \circ B}(x)=0$ and $\nu_{A \circ B}(x)=1$.
Theorem 3.12. If $A=\left\langle\mu_{A}, \nu_{A}\right\rangle$ and $B=\left\langle\mu_{B}, \nu_{B}\right\rangle$ are intuitionistic fuzzy ideals in a $\Gamma$ - ring $M$ then $A \circ B$ is an intuitionistic fuzzy ideal in $M$.
Proof. For any $x, y \in M$, we have

$$
\begin{aligned}
& \mu_{A \circ B}(x-y) \\
& =\bigvee\left\{\bigwedge_{1 \leq i \leq k} \mu_{A}\left(u_{i}\right) \wedge \mu_{B}\left(v_{i}\right): x-y=\sum_{1}^{k} u_{i} \alpha v_{i}, u_{i}, v_{i} \in M, \alpha \in \Gamma \text { and } k \in N\right\} \\
& \geq \bigvee\left\{\left(\bigwedge_{1 \leq i \leq m} \mu_{A}\left(a_{i}\right) \wedge \mu_{B}\left(b_{i}\right)\right) \wedge\left(\bigwedge_{1 \leq i \leq n} \mu_{A}\left(-c_{i}\right) \wedge \mu_{B}\left(d_{i}\right)\right)\right. \\
& \left.: x=\sum_{1}^{m} a_{i} \alpha b_{i},-y=\sum_{1}^{n}-c_{i} \alpha d_{i}, a_{i}, b_{i},-c_{i}, d_{i} \in M, \alpha \in \Gamma \text { and } m, n \in N\right\} \\
& =\bigvee\left\{\left(\bigwedge_{1 \leq i \leq m} \mu_{A}\left(a_{i}\right) \wedge \mu_{B}\left(b_{i}\right)\right) \wedge\left(\bigwedge_{1 \leq i \leq n} \mu_{A}\left(c_{i}\right) \wedge \mu_{B}\left(d_{i}\right)\right)\right. \\
& \left.: x=\sum_{1}^{m} a_{i} \alpha b_{i}, y=\sum_{1}^{n} c_{i} \alpha d_{i}, a_{i}, b_{i}, c_{i}, d_{i} \in M, \alpha \in \Gamma \text { and } m, n \in N\right\} \\
& =\bigvee\left\{\bigwedge_{1 \leq i \leq m} \mu_{A}\left(a_{i}\right) \wedge \mu_{B}\left(b_{i}\right): x=\sum_{1}^{m} a_{i} \alpha b_{i}, a_{i}, b_{i}, \in M, \alpha \in \Gamma \text { and } m \in N\right\} \\
& \wedge \bigvee\left\{\bigwedge_{1 \leq i \leq n} \mu_{A}\left(c_{i}\right) \wedge \mu_{B}\left(d_{i}\right): y=\sum_{1}^{n} c_{i} \alpha d_{i}, c_{i}, d_{i}, \in M, \alpha \in \Gamma \text { and } n \in N\right\} \\
& =\mu_{A \circ B}(x) \wedge \mu_{A \circ B}(y) \\
& \nu_{A \circ B}(x-y) \\
& =\bigwedge\left\{\underset{1 \leq i \leq k}{\bigvee} \nu_{A}\left(u_{i}\right) \vee \nu_{B}\left(v_{i}\right): x-y=\sum_{1}^{k} u_{i} \alpha v_{i}, u_{i}, v_{i} \in M, \alpha \in \Gamma \text { and } k \in N\right\} \\
& \leq \bigwedge\left\{\left(\underset{1 \leq i \leq m}{\bigvee} \nu_{A}\left(a_{i}\right) \vee \nu_{B}\left(b_{i}\right)\right) \vee\left(\underset{1 \leq i \leq n}{\bigvee} \nu_{A}\left(-c_{i}\right) \vee \nu_{B}\left(d_{i}\right)\right)\right. \\
& \left.: x=\sum_{1}^{m} a_{i} \alpha b_{i},-y=\sum_{1}^{n}-c_{i} \alpha d_{i}, a_{i}, b_{i},-c_{i}, d_{i} \in M, \alpha \in \Gamma \text { and } m, n \in N\right\} \\
& =\bigwedge\left\{\left(\underset{1 \leq i \leq m}{\bigvee} \nu_{A}\left(a_{i}\right) \vee \nu_{B}\left(b_{i}\right)\right) \vee\left(\underset{1 \leq i \leq n}{\bigvee} \nu_{A}\left(c_{i}\right) \vee \nu_{B}\left(d_{i}\right)\right)\right. \\
& \left.: x=\sum_{1}^{m} a_{i} \alpha b_{i}, y=\sum_{1}^{n} c_{i} \alpha d_{i}, a_{i}, b_{i}, c_{i}, d_{i} \in M, \alpha \in \Gamma \text { and } m, n \in N\right\} \\
& =\bigwedge\left\{\underset{1 \leq i \leq m}{\bigvee} \nu_{A}\left(a_{i}\right) \vee \nu_{B}\left(b_{i}\right): x=\sum_{1}^{m} a_{i} \alpha b_{i}, a_{i}, b_{i} \in M, \alpha \in \Gamma \text { and } m \in N\right\} \\
& \vee \bigwedge\left\{\underset{1 \leq i \leq n}{\bigvee} \nu_{A}\left(c_{i}\right) \vee \nu_{B}\left(d_{i}\right): y=\sum_{1}^{n} c_{i} \alpha d_{i}, c_{i}, d_{i} \in M, \alpha \in \Gamma \text { and } n \in N\right\} \\
& =\nu_{A \circ B}(x) \vee \nu_{A \circ B}(y) .
\end{aligned}
$$

Also, we have

$$
\begin{aligned}
& \mu_{A \circ B}(x)=\bigvee\left\{\bigwedge_{1 \leq i \leq m} \mu_{A}\left(a_{i}\right) \wedge \mu_{B}\left(b_{i}\right): x=\sum_{1}^{m} a_{i} \alpha b_{i}, a_{i}, b_{i} \in M, \alpha \in \Gamma \text { and } m \in N\right\} \\
& \leq \bigvee\left\{\bigwedge_{1 \leq i \leq m} \mu_{A}\left(a_{i}\right) \wedge \mu_{B}\left(b_{i} \alpha y\right): x \alpha y=\sum_{1}^{m} a_{i} \alpha\left(b_{i} \alpha y\right), a_{i}, b_{i} \alpha y \in M, \alpha \in \Gamma \text { and } m \in\right. \\
& N\} \\
& \leq \bigvee\left\{\bigwedge_{1 \leq i \leq m} \mu_{A}\left(u_{i}\right) \wedge \mu_{B}\left(v_{i}\right): x \alpha y=\sum_{1}^{m} u_{i} \alpha v_{i}, u_{i}, v_{i} \in M, \alpha \in \Gamma \text { and } m \in N\right\} \\
& =\mu_{A \circ B}(x \alpha y) \\
& \nu_{A \circ B}(x)=\bigwedge\left\{\bigvee_{1 \leq i \leq m}^{\bigvee} \nu_{A}\left(a_{i}\right) \vee \nu_{B}\left(b_{i}\right): x=\sum_{1}^{m} a_{i} \alpha b_{i}, a_{i}, b_{i} \in M, \alpha \in \Gamma \text { and } m \in N\right\} \\
& \geq \bigwedge\left\{\bigvee_{1 \leq i \leq m} \nu_{A}\left(a_{i}\right) \vee \nu_{B}\left(b_{i} \alpha y\right): x \alpha y=\sum_{1}^{m} a_{i} \alpha\left(b_{i} \alpha y\right), a_{i}, b_{i} \alpha y \in M, \alpha \in \Gamma \text { and } m \in\right. \\
& N\} \\
& \geq \bigwedge\left\{\bigvee_{1 \leq i \leq m} \nu_{A}\left(u_{i}\right) \vee \nu_{B}\left(v_{i}\right): x \alpha y=\sum_{1}^{m} u_{i} \alpha v_{i}, u_{i}, v_{i} \in M, \alpha \in \Gamma \text { and } m \in N\right\} \\
& =\nu_{A \circ B}(x \alpha y) .
\end{aligned}
$$

Hence $\mu_{A \circ B}(x \alpha y) \geq \mu_{A \circ B}(x)$ and $\nu_{A \circ B}(x \alpha y) \leq \nu_{A \circ B}(x)$.
Similarly we get $\mu_{A \circ B}(x \alpha y) \geq \mu_{A \circ B}(y)$ and $\nu_{A \circ B}(x \alpha y) \leq \nu_{A \circ B}(y)$.
Therefore $A \circ B=\left(\mu_{A \circ B}, \nu_{A \circ B}\right)$ is an intuitionistic fuzzy ideal of $M$.
Definition 3.13[2]. A function $f: M \rightarrow N$, where $M$ and $N$ are $\Gamma$ - rings, is said to be a $\Gamma$ - homomorphism if $f(a+b)=f(a)+f(b), f(a \alpha b)=f(a) \alpha f(b)$, for all $a, b \in M$ and $\alpha \in \Gamma$.
Definition 3.14[2]. A function $f: M \rightarrow N$, where $f$ is a $\Gamma$ - homomorphism and $M$ and $N$ are $\Gamma$ - rings, is said to be a $\Gamma$ - endomorphism if $N \subseteq M$.

Definition 3.15[2]. Let $f: X \rightarrow Y$ be a mapping of $\Gamma$ - rings and $A$ be an intuitionistic fuzzy set of $Y$. Then the map $f^{-1}(A)$ is the pre-image of $A$ under $f$, if $\mu_{f^{-1}(A)}(X)=\mu_{A}(f(x))$ and $\nu_{f^{-1}(A)}(X)=\nu_{A}(f(x))$, for all $x \in X$.

Theorem 3.16 Let $f$ be a $\Gamma$ - homomorphism of $M$. If the IFS $A=\left\langle\mu_{A}, \nu_{A}\right\rangle$ is an intuitionistic fuzzy left (resp. right) ideal of $M$, then $B=\left\langle\mu_{f^{-1}(A)}, \nu_{f^{-1}(A)}\right\rangle$ is an intuitionistic fuzzy left (resp. right) ideal of $M$.
Proof. For any $x, y \in M, \alpha \in \Gamma$, we have

$$
\begin{aligned}
& \qquad \mu_{f^{-1}(A)}(x-y)=\mu_{A}(f(x-y)) \\
& \quad=\mu_{A}(f(x)-f(y)) \geq \mu_{A}(f(x)) \wedge \mu_{A}(f(y))=\mu_{f^{-1}(A)}(x) \wedge \mu_{f^{-1}(A)}(y) \\
& \text { and } \mu_{f-1}(A)(x \alpha y)=\mu_{A}(f(x \alpha y))=\mu_{A}(f(x) \alpha f(y)) \geq \mu_{f^{-1}(A)}(y) \\
& \text { Similarly, } \quad \nu_{f-1}(A)(x-y)=\nu_{A}(f(x-y)) \\
& \quad=\nu_{A}(f(x)-f(y)) \leq \nu_{A}(f(x)) \vee \nu_{A}(f(y))=\nu_{f^{-1}(A)}(x) \vee \nu_{f^{-1}(A)}(y) \\
& \text { and } \quad \nu_{f^{-1}(A)}(x \alpha y)=\nu_{A}(f(x \alpha y))=\nu_{A}(f(x) \alpha f(y)) \leq \nu_{f^{-1}(A)}(y) \text {. }
\end{aligned}
$$

Hence $B$ is an intuitionistic fuzzy left (resp. right) ideal of $M$.
Theorem 3.17. If $A=\left\langle\mu_{A}, \nu_{A}\right\rangle$ is an intuitionistic fuzzy set in $M$ such that the
non-empty sets $U\left(\mu_{A} ; t\right)$ and $L\left(\nu_{A} ; t\right)$ are ideals of $M$ for all $t \in[0,1]$, then $A$ is an intuitionistic fuzzy left (resp. right) ideal of $M$.
Proof. Suppose that there exists $x_{0}, y_{0} \in M$ such that
$\mu_{A}\left(x_{0}-y_{0}\right)<\left(\mu_{A}\left(x_{0}\right) \wedge \mu_{A}\left(y_{0}\right)\right)$.
Let $t_{0}=\frac{1}{2}\left\{\mu_{A}\left(x_{0}-y_{0}\right)+\left(\mu_{A}\left(x_{0}\right) \wedge \mu_{A}\left(y_{0}\right)\right)\right\}$.
Then $\left(\mu_{A}\left(x_{0}\right) \wedge \mu_{A}\left(y_{0}\right)\right) \geq t_{0}>\mu_{A}\left(x_{0}-y_{0}\right)$.
It follows that $x_{o}, y_{o} \in U\left(\mu_{A} ; t_{0}\right)$ and $x_{0}-y_{0} \notin U\left(\mu_{A} ; t_{0}\right)$.
This is a contradiction.
Hence $\mu_{A}(x-y) \geq\left(\mu_{A}(x) \wedge \mu_{A}(y)\right)$, for all $x, y \in M$.
Now let $x_{0}, y_{0} \in M$ and $\alpha \in \Gamma$ such that $\mu_{A}\left(x_{0} \alpha y_{0}\right)<\mu_{A}\left(y_{0}\right)$.
Let $t_{0}=\frac{1}{2}\left\{\mu_{A}\left(x_{0} \alpha y_{0}\right)+\mu_{A}\left(y_{0}\right)\right\}$.
Then we get $\mu_{A}\left(x_{0} \alpha y_{0}\right) \leq t_{0}<\mu_{A}\left(y_{0}\right)$.
It follows that $y_{0} \in U\left(\mu_{A} ; t_{0}\right)$ and $x_{0} \alpha y_{0} \notin U\left(\mu_{A} ; t_{0}\right)$.
This is a contradiction.
Thus $\mu_{A}\left(x_{0} \alpha y_{0}\right) \geq \mu_{A}\left(y_{0}\right) \quad$ (resp. $\left.\mu_{A}\left(x_{0} \alpha y_{0}\right) \geq \mu_{A}\left(x_{0}\right)\right)$.
Similarly, suppose that there exists $x_{0}, y_{0} \in M$ such that
$\nu_{A}\left(x_{0}-y_{0}\right)>\left\{\nu_{A}\left(x_{0}\right) \vee \nu_{A}\left(y_{0}\right)\right\}$.
Let $t_{0}=\frac{1}{2}\left\{\nu_{A}\left(x_{0}-y_{0}\right)+\left(\nu_{A}\left(x_{0}\right) \vee \nu_{A}\left(y_{0}\right)\right)\right\}$.
Then $\left(\nu_{A}\left(x_{0}\right) \vee \nu_{A}\left(y_{0}\right)\right) \leq t_{0}<\nu_{A}\left(x_{0}-y_{0}\right)$.
It follows that $x_{o}, y_{o} \in L\left(\mu_{A} ; t_{0}\right)$ and $x_{0}-y_{0} \notin L\left(\mu_{A} ; t_{0}\right)$.
This is a contradiction.
Hence $\nu_{A}(x-y) \leq\left(\nu_{A}(x) \vee \nu_{A}(y)\right)$, for all $x, y \in M$.
Now let $x_{0}, y_{0} \in M$ and $\alpha \in \Gamma$ such that $\nu_{A}\left(x_{0} \alpha y_{0}\right)>\nu_{A}\left(y_{0}\right)$.
Let $t_{0}=\frac{1}{2}\left\{\nu_{A}\left(x_{0} \alpha y_{0}\right)+\nu_{A}\left(y_{0}\right)\right\}$.
Then we get $\nu_{A}\left(x_{0} \alpha y_{0}\right)>t_{0}>\nu_{A}\left(y_{0}\right)$.
It follows that $y_{0} \in L\left(\mu_{A} ; t_{0}\right)$ and $x_{0} \alpha y_{0} \notin L\left(\nu_{A} ; t_{0}\right)$.
This is a contradiction.
Thus $\nu_{A}\left(x_{0} \alpha y_{0}\right) \leq \nu_{A}\left(y_{0}\right) \quad\left(\right.$ resp. $\left.\nu_{A}\left(x_{0} \alpha y_{0}\right) \leq \nu_{A}\left(x_{0}\right)\right)$.
Hence $A$ is an intuitionistic fuzzy left (resp. right) ideal of $M$.

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# A new class of analytic functions based on Ruscheweyh derivative 

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AbSTRACT: In this paper we introduce a new class $Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$ consisting of analytic functions with negative coefficients and investigate various properties and characterization of the class. The results include coefficient estimates, distortion theorem, closure theorems and integral operators for the class $Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$. Also radii of close-to-convexity, starlikeness and convexity are determined

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## 1 Introduction

Let $\mathcal{A}$ denote the class of functions of the form

$$
\begin{equation*}
f(z)=z+\sum_{n=2}^{\infty} a_{n} z^{n} \tag{1.1}
\end{equation*}
$$

that are analytic in the open unit disc $\Delta:=\{z \in \mathbb{C}:|z|<1\}$. Let $\mathcal{S}$ be a subclass of $\mathcal{A}$ consisting of functions univalent in $\Delta$.

Let $f$ and $g$ be functions analytic in $\Delta$. Then we say that $f$ is subordinate to $g$ if there exists a Schwarz function $w(z)$, analytic in $\Delta$ with $w(0)=0 \quad$ and $\quad|w(z)|<$ $1(z \in \Delta)$, such that

$$
f(z)=g(\omega(z)) \quad(z \in \Delta)
$$

We denote this subordination by

$$
f \prec g \quad \text { or } \quad f(z) \prec g(z) \quad(z \in \Delta) .
$$

In particular, if $g$ is univalent in $\Delta$, the above subordination is equivalent to

$$
f(0)=g(0) \quad \text { and } \quad f(\Delta) \subset g(\Delta)
$$

For functions $f \in \mathcal{A}$ given by (1.1) and $g \in \mathcal{A}$ given by $g(z)=z+\sum_{n=2}^{\infty} b_{n} z^{n}$, we define the Hadamard product (or Convolution) of $f$ and $g$ by

$$
\begin{equation*}
(f * g)(z)=z+\sum_{n=2}^{\infty} a_{n} b_{n} z^{n}, \quad z \in \Delta \tag{1.2}
\end{equation*}
$$

By using the Hadamard product, Ruscheweyh [10] defined an operator

$$
\begin{equation*}
D^{\gamma} f(z)=\frac{z}{(1-z)^{\gamma+1}} * f(z), \quad \gamma \geq-1 \tag{1.3}
\end{equation*}
$$

Ruscheweyh [10] observed that

$$
\begin{equation*}
D^{n} f(z)=\frac{z\left(z^{n-1} f(z)\right)^{(n)}}{n!} \tag{1.4}
\end{equation*}
$$

where $n=\gamma \in \mathbb{N}_{0}=\{0,1,2, \ldots\}$. This symbol $D^{n} f(z), n \in \mathbb{N}_{0}$ is called by Al-Amiri [1], the $n^{\text {th }}$ order Ruscheweyh derivative of $f(z)$.

We note that $D^{0} f(z)=f(z), D^{1} f(z)=z f^{\prime}(z)$, and

$$
\begin{equation*}
D^{n} f(z)=z+\sum_{k=2}^{\infty} \sigma(n, k) a_{k} z^{k} \tag{1.5}
\end{equation*}
$$

where

$$
\begin{equation*}
\sigma(n, k)=\binom{n+k-1}{n} \tag{1.6}
\end{equation*}
$$

Let $T$ denote the subclass of $S$ consisting of functions of the form

$$
\begin{equation*}
f(z)=z-\sum_{k=2}^{\infty} a_{k} z^{k}, \quad a_{k} \geq 0 \tag{1.7}
\end{equation*}
$$

Several new classes of analytic functions defined by Ruscheweyh derivatives have been studied and continue to be introduced and investigated in the literature (See for example $[12,13,8]$ to mention a few interesting studies). Several investigations on functions with negative coefficients have been done. (See for example $[3,2,6,9,11]$ ). In particular, results on functions with negative coefficients related to Ruscheweyh derivatives have been derived. See for example, $[4,5,7]$ ).

Motivated by the aforementioned works, we introduce a new class $Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$ by using $m^{t h}$ and $n^{t h}$ order Ruscheweyh derivative of $f(z)$.

Definition 1 Let $\Phi(z)=z+\sum_{k=2}^{\infty} \beta_{k} z^{k}$, and $\Psi(z)=z+\sum_{k=2}^{\infty} \gamma_{k} z^{k}$ be fixed analytic functions in $\Delta$ and $\beta_{k}>0, \gamma_{k}>0, k \geq 2$. We define a class $Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$
consisting of analytic functions of the form

$$
f(z)=z-\sum_{k=2}^{\infty} a_{k} z^{k}, a_{k} \geq 0
$$

which satisfy the subordination condition

$$
\begin{equation*}
(1-\lambda)\left(D^{m}(f * \Phi)(z)\right)^{\prime}+\lambda\left(D^{n}(f * \Psi)(z)\right)^{\prime} \prec \frac{1+A z}{1+B z} \tag{1.8}
\end{equation*}
$$

for $z \in \Delta$, where $\lambda \geq 0,-1 \leq A<B \leq 1,0<B \leq 1$ and $m, n \in N_{0}$.
By specializing the parameters $m, n, \lambda, A$ and $B$, and the functions $\Phi$ and $\Psi$, we obtain the subclasses studied by various authors as listed below:
(i) $Q_{m}^{m+1}\left(\frac{z}{1-z}, \frac{z}{1-z}, \lambda, A, B\right)=Q(m, \lambda, A, B)[4]$
(ii) $Q_{0}^{1}\left(\frac{z}{1-z}, \frac{z}{1-z}, \lambda, 2 \alpha-1,1\right)=R(\lambda, \alpha)(0 \leq \alpha<1)[3]$
(iii) $Q_{0}^{1}\left(\frac{z}{1-z}, \frac{z}{1-z}, 0,2 \alpha-1,1\right)=T^{*}(\alpha)(0 \leq \alpha<1)[11,2]$
(iv) $Q_{m}^{m+1}\left(\frac{z}{1-z}, \frac{z}{1-z}, 0,2 \alpha-1,1\right)=Q_{n}(\alpha)(0 \leq \alpha<1)[14]$
(v) $Q_{0}^{1}\left(\frac{z}{1-z}, \frac{z}{1-z}, 0,(2 \alpha-1) \beta, \beta\right)=p^{*}(\alpha, \beta)$
$(0 \leq \alpha<1,0<\beta \leq 1)[6]$
(vi) $Q_{0}^{1}\left(\frac{z}{1-z}, \frac{z}{1-z}, 0,((1+\mu) \alpha-1) \beta, \mu \beta\right)=p^{*}(\alpha, \beta, \mu)$ $(0 \leq \alpha<1,0<\beta \leq 1,0 \leq \mu<1)[9]$
It is of interest to note that the class $Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$ gives several well-known subclasses of functions for suitable choices of $\Phi(z), \Psi(z)$ and $m, n$. We obtain coefficient inequality, coefficient estimate, distortion theorem, extreme points and integral representation for functions in the class $Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$.

## 2 Coefficient Estimates

Theorem 1 A function $f(z)=z-\sum_{k=2}^{\infty} a_{k} z^{k} \in Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$ if and only if

$$
\begin{equation*}
\sum_{k=2}^{\infty} k\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right] a_{k} \leq \frac{B-A}{B+1} \tag{2.1}
\end{equation*}
$$

Proof. Suppose $f(z) \in Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$. In view of the definition of subordination we get,

$$
\begin{equation*}
h(z)=(1-\lambda)\left(D^{m}(f * \Phi)(z)\right)^{\prime}+\lambda\left(D^{n}(f * \Psi)(z)\right)^{\prime}=\frac{1+A w(z)}{1+B w(z)} \tag{2.2}
\end{equation*}
$$

$-1 \leq A<B \leq 1,0<B \leq 1, z \in \Delta$, and $|\omega(z)|<1$. From (2.2), we get

$$
w(z)=\frac{1-h(z)}{B h(z)-A}
$$

In view of (2.2), one can easily obtain through a simple computation that,

$$
h(z)=1-\sum_{k=0}^{\infty} k\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right] a_{k} z^{k-1}
$$

and $|w(z)|<1$ implies

$$
\begin{equation*}
\left|\frac{\sum_{k=2}^{\infty} k\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right] a_{k} z^{k-1}}{(B-A)-B \sum_{k=2}^{\infty} k\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right] a_{k} z^{k-1}}\right|<1 \tag{2.3}
\end{equation*}
$$

and hence,(2.1)holds.
Conversely, Suppose $f \in T$ and satisfies (2.1). For $|z|=r, 0 \leq r<1$, we have by (2.1),

$$
\begin{array}{r}
|1-h(z)|-|B h(z)-A| \leq \sum_{k=2}^{\infty} k\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right] a_{k} r^{k-1} \\
-(B-A)+B \sum_{k=2}^{\infty} k\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right] a_{k} r^{k-1} \leq 0
\end{array}
$$

which gives (2.2) and hence follows that

$$
(1-\lambda)\left(D^{m}(f * \Phi)(z)\right)^{\prime}+\lambda\left(D^{n}(f * \Psi)(z)\right)^{\prime}=\frac{1+A w(z)}{1+B w(z)}
$$

$z \in \Delta,-1 \leq A<B \leq 1,0<B \leq 1$. Hence, $f(z) \in Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$.
Finally the function $f(z)$ given by

$$
\begin{equation*}
f(z)=z-\frac{B-A}{k(1+B)\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right]} z^{k}, \quad k \geq 2 \tag{2.4}
\end{equation*}
$$

is an extremal function for the class.

Corollary 1 Let the function $f(z)$ defined by (1.7) be in the class $Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$. Then we have

$$
\begin{equation*}
a_{n} \leq \frac{B-A}{k(1+B)\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right]}, \quad k \geq 2 \tag{2.5}
\end{equation*}
$$

Theorem 2

$$
Q_{m}^{n}\left(\Phi, \Psi, \lambda_{2}, A, B\right) \subseteq Q_{m}^{n}\left(\Phi, \Psi, \lambda_{1}, A, B\right)
$$

for $-1 \leq A<B \leq 1,0<B \leq 1, \lambda_{2} \geq \lambda_{1} \geq 0$ and $n>m$.
Proof. Let $f(z) \in Q_{m}^{n}\left(\Phi, \Psi, \lambda_{2}, A, B\right)$.
$\sum_{k=2}^{\infty} k\left[\left(1-\lambda_{1}\right) \sigma(m, k) \beta_{k}+\lambda_{1} \sigma(n, k) \gamma_{k}\right] a_{k}$

$$
\leq \sum_{k=2}^{\infty} k\left[\sigma(m, k) \beta_{k}+\left(\sigma(n, k) \gamma_{k}-\sigma(m, k) \beta_{k}\right) \lambda_{2}\right] a_{k} \leq \frac{B-A}{1+B}
$$

Therefore $f(z) \in Q_{m}^{n}\left(\Phi, \Psi, \lambda_{1}, A, B\right)$. Hence the proof of Theorem 2 is complete.

## 3 Closure Theorems

Let the functions $f_{i}(z)$ be defined, for $i=1,2, \ldots, \ell$ by

$$
\begin{equation*}
f_{i}(z)=z-\sum_{k=2}^{\infty} a_{k, \ell} z^{k}, \quad a_{k, \ell} \geq 0 \tag{3.1}
\end{equation*}
$$

We shall prove the following results for the closure of functions in the class $Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$.

Theorem 3 Let the function $f_{i}(z)$ defined by (3.1) be in the class $Q_{m}^{n}\left(\Phi, \Psi, \lambda, A_{i}, B_{i}\right)$, for $i=1,2, \ldots, \ell$. Then the function $h(z)$ defined by

$$
\begin{equation*}
h(z)=z-\frac{1}{\ell} \sum_{k=2}^{\infty}\left(\sum_{i=1}^{\ell} a_{k, i} z^{k}\right) \tag{3.2}
\end{equation*}
$$

is in the class $Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$,

$$
\begin{equation*}
\text { where } A=\min _{1 \leq i \leq \ell}\left\{A_{i}\right\} \quad \text { and } \quad B=\max _{1 \leq i \leq \ell}\left\{B_{i}\right\} \tag{3.3}
\end{equation*}
$$

Proof. Since $f_{i}(z) \in Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$, for $i=1,2, \ldots, \ell$, we have

$$
\sum_{k=2}^{\infty} k\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right] a_{k, i} \leq \frac{B_{i}-A_{i}}{1+B_{i}}, \text { by Theorem } 1
$$

Hence we obtain,

$$
\begin{aligned}
\sum_{k=2}^{\infty} k[(1-\lambda) \sigma(m, & \left.k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right]\left[\frac{1}{\ell} \sum_{i=1}^{\ell} a_{k, i}\right] \\
& \leq \frac{1}{\ell} \sum_{i=1}^{\ell}\left\{\sum_{k=2}^{\infty} k\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right] a_{k, i}\right\} \\
& \leq \frac{1}{\ell} \sum_{i=1}^{\ell} \frac{B_{i}-A_{i}}{1+B_{i}} \leq \frac{B-A}{1+B}
\end{aligned}
$$

Thus, we get

$$
\sum_{k=2}^{\infty} k\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right]\left[\frac{1}{\ell} \sum_{i=1}^{\ell} a_{k, i}\right] \leq \frac{B-A}{1+B}
$$

given (3.3), which shows that $f(z) \in Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$.
Theorem 4 Let the functions $f_{i}(z)(i=1,2, \ldots, \ell)$ defined by (3.1) be in the class $Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$. Then the function $h(z)$ defined by $h(z)=\sum_{i=1}^{\ell} d_{i} f_{i}(z)$ is also in the same class $Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$, where $\sum_{i=1}^{\ell} d_{i}=1$.
Proof. According to definition of $h(z)$ we can write that

$$
h(z)=z-\sum_{k=2}^{\infty}\left(\sum_{i=1}^{\ell} d_{i} a_{k, i}\right) z^{k} .
$$

Further, since $f_{i}(z)$ are in $Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$ for every $i=1,2, \ldots, \ell$, we get

$$
\sum_{k=2}^{\infty} k\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right] a_{k, i} \leq \frac{B-A}{1+B}
$$

for every $i=1,2, \ldots, \ell$.
Hence, we can see that

$$
\begin{aligned}
\sum_{k=2}^{\infty} k[(1-\lambda) \sigma(m & \left., k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right]\left(\sum_{i=1}^{\ell} d_{i} a_{k, i}\right) \\
& =\sum_{i=1}^{\ell} d_{i}\left(\sum_{k=2}^{\infty} k\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right] a_{k, i}\right) \\
& \leq\left(\sum_{i=1}^{\ell} d_{i}\right) \frac{B-A}{1+B}=\frac{B-A}{1+B}
\end{aligned}
$$

This proves that the function $h(z) \in Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$. Thus we have the theorem.

Corollary 2 The class $Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$ is closed under convex linear combination.
Proof. Putting $\ell=2$ in the above theorem, we prove the corollary.
Theorem 5 Let

$$
f_{1}(z)=z
$$

and

$$
f_{k}(z)=z-\frac{B-A}{k(B+1)\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right]} z^{k}, k \geq 2
$$

Then $f(z) \in Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$ if and only if it can be expressed in the form

$$
\begin{equation*}
f(z)=\sum_{k=1}^{\infty} \mu_{k} f_{k}(z) \tag{3.4}
\end{equation*}
$$

where $\mu_{k} \geq 0(k \geq 1)$ and $\sum_{k=1}^{\infty} \mu_{k}=1$.
Proof. Assume that $f(z)=\sum_{k=1}^{\infty} \mu_{k} f_{k}(z)$. Then

$$
\begin{equation*}
f(z)=z-\sum_{k=2}^{\infty} \frac{(B-A) \mu_{k}}{k(B+1)\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right]} z^{k} \tag{3.5}
\end{equation*}
$$

Then it follows that

$$
\begin{aligned}
& \sum_{k=2}^{\infty} \frac{k(B+1)\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right]}{B-A} \\
& \cdot \frac{(B-A) \mu_{k}}{k(B+1)\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right]}=\sum_{k=2}^{\infty} \mu_{k} \\
&=1-\mu_{1} \\
& \leq 1
\end{aligned}
$$

Then $f(z) \in Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$.
Conversely, assume that $f(z)$ defined by (1.7) belongs to the class $Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$. Then

$$
a_{k} \leq \frac{(B-A)}{k(B+1)\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right]} \quad(k \geq 2)
$$

Setting,

$$
\mu_{k}=\frac{k(B+1)\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right]}{(B-A)}, \quad \mu_{1} \quad=1-\sum_{k=2}^{\infty} \mu_{k}
$$

we can see that $f(z)$ can be expressed in the form (3.4). This completes the proof.

## 4 A Set of Distortion Inequalities

Theorem 6 Let the function $f(z)$ defined by (1.7) be in the class $Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$ and let

$$
\left[(1-\lambda) \sigma(m, 2) \beta_{2}+\lambda \sigma(n, 2) \gamma_{2}\right] \leq\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right]
$$

Then, we have for $|z|=r<1$,

$$
\begin{equation*}
r-\frac{B-A}{2 \delta_{2}} r^{2} \leq|f(z)| \leq r+\frac{B-A}{2 \delta_{2}} r^{2} \tag{4.1}
\end{equation*}
$$

and

$$
\begin{equation*}
1-\frac{B-A}{\delta_{2}} r \leq\left|f^{\prime}(z)\right| \leq 1+\frac{B-A}{\delta_{2}} r \tag{4.2}
\end{equation*}
$$

The results are sharp for the function

$$
\begin{equation*}
f(z)=z-\frac{B-A}{2 \delta_{2}} r^{2} \tag{4.3}
\end{equation*}
$$

where $\delta_{2}=2(1+B)\left[(1-\lambda) \sigma(m, 2) \beta_{2}+\lambda \sigma(n, 2) \gamma_{2}\right]$.
Proof. Since $k\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right]$ is an increasing function of $k(k \geq$ $2)$, and $f(z) \in Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$, by Theorem 1, we have

$$
\begin{aligned}
2\left[(1-\lambda) \sigma(m, 2) \beta_{2}+\lambda \sigma(n, 2) \gamma_{2}\right] \sum_{k=2}^{\infty} a_{k} & \leq \sum_{k=2}^{\infty} k\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right] a_{k} \\
& \leq \frac{B-A}{1+B}
\end{aligned}
$$

That is

$$
\begin{aligned}
\sum_{k=2}^{\infty} a_{k} & \leq \frac{B-A}{2(1+B)\left[(1-\lambda) \sigma(m, 2) \beta_{2}+\lambda \sigma(n, 2) \gamma_{2}\right]}=\frac{B-A}{2 \delta_{2}} \\
|f(z)| & \leq z+\sum_{k=2}^{\infty}\left|a_{k}\right||z|^{k} \\
& \leq r+\sum_{k=2}^{\infty} a_{k} r^{k} \leq r+r^{2} \sum_{k=2}^{\infty} a_{k} \\
& \leq r+r^{2}\left[\frac{B-A}{2(1+B)\left[(1-\lambda) \sigma(m, 2) \beta_{2}+\lambda \sigma(n, 2) \gamma_{2}\right]}\right] \\
& =r+r^{2}\left(\frac{B-A}{2 \delta_{2}}\right) \text { and } \\
|f(z)| & \geq r-\sum_{k=2}^{\infty} a_{k} r^{k} \geq r-r^{2} \sum_{k=2}^{\infty} a_{k}=r-r^{2}\left(\frac{B-A}{2 \delta_{2}}\right)
\end{aligned}
$$

Hence (4.1) follows.
Also, in view of the inequality (2.1), we have

$$
\left[(1-\lambda) \sigma(m, 2) \beta_{2}+\lambda \sigma(n, 2) \gamma_{2}\right] \sum_{k=2}^{\infty} k a_{k} \leq \frac{B-A}{1+B}
$$

which gives, $\sum_{k=2}^{\infty} k a_{k} \leq \frac{B-A}{\delta_{2}}$. Thus

$$
\left|f^{\prime}(z)\right| \leq 1+\sum_{k=2}^{\infty} k a_{k} r^{k-1} \leq 1+r \sum_{k=2}^{\infty} k a_{k} \leq 1+\frac{B-A}{\delta_{2}} r
$$

Similarly we can prove the other inequality $\left|f^{\prime}(z)\right| \geq 1-\frac{B-A}{\delta_{2}} r$.
Hence (4.2) follows also.

## 5 Integral Operator

Theorem 7 Let the function $f(z)$ defined by (1.7) be in the class $Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$, and let $d$ be a real number such that $d>-1$. Then the function $F(z)$ defined by

$$
\begin{equation*}
F(z)=\frac{d+1}{z^{d}} \int_{0}^{z} t^{d-1} f(t) d t \tag{5.1}
\end{equation*}
$$

also belongs to the class $Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$.
Proof. From the representation of $F(z)$, it follows that

$$
\begin{equation*}
F(z)=z-\sum_{k=2}^{\infty}\left(\frac{d+1}{d+k}\right) a_{k} z^{k} \tag{5.2}
\end{equation*}
$$

Therefore,

$$
\begin{aligned}
& \sum_{k=2}^{\infty} k\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right]\left(\frac{d+1}{d+k}\right) a_{k} \\
& \leq \sum_{k=2}^{\infty} k\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right] a_{k} \leq \frac{B-A}{1+B}
\end{aligned}
$$

Hence, $F \in Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$.
Theorem 8 Let the function $F(z)=z-\sum_{k=2}^{\infty} a_{k} z^{k}\left(a_{k} \geq 0\right)$ be in the class
$Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$ and let $d$ be a real number such that $d>-1$. Then the function $f(z)$ defined by (5.1) is univalent in $|z|<R^{*}$, where

$$
\begin{equation*}
R^{*}=\inf _{k}\left[\frac{(1+B)\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right](d+1)}{(B-A)(d+m)}\right]^{1 / k-1} \quad(k \geq 2) \tag{5.3}
\end{equation*}
$$

The result is sharp.

Proof. From (5.1), we have

$$
\begin{aligned}
f(z) & =\frac{z^{1-d}\left(z^{d} F(z)\right)^{\prime}}{d+1} \quad(d>-1) \\
& =z-\sum_{k=2}^{\infty}\left(\frac{d+k}{d+1}\right) a_{k} z^{k}
\end{aligned}
$$

In order to obtain the required result it suffices to show that

$$
\begin{equation*}
\left|f^{\prime}(z)-1\right|<1 \tag{5.4}
\end{equation*}
$$

for $|z|<R^{*}$. Now

$$
\left|f^{\prime}(z)-1\right| \leq \sum_{k=2}^{\infty} \frac{k(d+k)}{(d+1)} a_{k}|z|^{k-1}
$$

Thus $\left|f^{\prime}(z)-1\right|<1$, if $\sum_{k=2}^{\infty} \frac{k(d+k)}{(d+1)} a_{k}|z|^{k-1}<1$.
Hence by Theorem 1, (5.4) is true if

$$
\frac{k(d+k)}{(d+1)}|z|^{k-1} \leq \frac{k(B+1)\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right]}{B-A}
$$

or if

$$
|z|<\left[\frac{(B+1)\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right](d+1)}{(B-A)(d+k)}\right]^{1 / k-1}, \quad k \geq 2
$$

which proves that $f$ is univalent in $|z|<R^{*}$.
Sharpness follows if we take,

$$
f(z)=z-\frac{(B-A)(d+k)}{k(B+1)\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right](d+1)} z^{k} \quad k \geq 2
$$

## 6 Radii of Close-to-convexity, Starlikeness and Convexity

Theorem 9 Let the function $f(z)$ defined by (1.7) be in the class
$Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$. Then $f(z)$ is close-to-convex of order $\rho(0 \leq \rho<1)$ in $|z|<r_{1} \equiv$ $r_{1}(m, n, \Phi, \Psi, \lambda, A, \rho)$, where

$$
\begin{equation*}
r_{1}=\inf _{k}\left[\frac{(1-\rho)(1+B)\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right]}{(B-A)}\right]^{1 / k-1} \quad k \geq 2 \tag{6.1}
\end{equation*}
$$

The result is sharp with the extremal function $f(z)$ is given by (2.4).

Proof. We must show that $\left|f^{\prime}(z)-1\right| \leq 1-\rho$ for $|z|<r_{1}$. We have

$$
\left|f^{\prime}(z)-1\right| \leq \sum_{k=2}^{\infty} k a_{k}|z|^{k-1}
$$

Hence, $\left|f^{\prime}(z)-1\right| \leq 1-\rho$ if

$$
\begin{equation*}
\sum_{k=2}^{\infty}\left(\frac{k}{1-\rho}\right) a_{k}|z|^{k-1} \leq 1 \tag{6.2}
\end{equation*}
$$

By Theorem 1, we have

$$
\sum_{k=2}^{\infty} \frac{k(B+1)\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right] a_{k}}{(B-A)} \leq 1
$$

Hence, (6.2) will be true if

$$
\frac{k|z|^{k-1}}{1-\rho} \leq \frac{k(B+1)\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right] a_{k}}{(B-A)}
$$

or if

$$
\begin{equation*}
|z| \leq\left[\frac{(1-\rho)(B+1)\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right] a_{k}}{(B-A)}\right]^{1 / k-1}, \quad k \geq 2 \tag{6.3}
\end{equation*}
$$

The theorem follows easily from (6.3).
Theorem 10 Let the function $f(z)$ defined by (1.7) be in the class $Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$, then $f(z)$ is starlike of order $\rho(0 \leq \rho<1)$ in $|z|<r_{2} \equiv r_{2}(m, n, \Phi, \Psi, \lambda, A, B, \rho)$ where

$$
\begin{equation*}
r_{2}=\inf _{k}\left[\frac{(1-\rho) k(1+B)\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right]}{(B-A)(k-\rho)}\right]^{1 / k-1} \quad(k \geq 2) \tag{6.4}
\end{equation*}
$$

The result is sharp with the extremal function $f(z)$ given by (2.4).
Proof. It is sufficient to show $\left|\frac{z f^{\prime}(z)}{f(z)}-1\right| \leq 1-\rho$ for $|z|<r_{2}$. By making use of the inequality,

$$
\begin{equation*}
\left|\frac{z f^{\prime}(z)}{f(z)}-1\right| \leq \frac{\sum_{k=2}^{\infty}(k-1) a_{k}|z|^{k-1}}{1-\sum_{k=2}^{\infty} a_{k}|z|^{k-1}} \tag{6.5}
\end{equation*}
$$

we get, $\left|\frac{z f^{\prime}(z)}{f(z)}-1\right| \leq 1-\rho$ if

$$
\begin{equation*}
\sum_{k=2}^{\infty} \frac{(k-\rho) a_{k}|z|^{k-1}}{1-\rho} \leq 1 \tag{6.6}
\end{equation*}
$$

Hence, by using (6.3), (6.6) will be true if

$$
\frac{(k-\rho)|z|^{k-1}}{1-\rho} \leq \frac{k(1+B)\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right]}{(B-A)}
$$

or if

$$
\begin{equation*}
|z| \leq\left[\frac{(1-\rho) k(1+B)\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right]}{(k-\rho)(B-A)}\right]^{1 / k-1}, \quad k \geq 2 \tag{6.7}
\end{equation*}
$$

The theorem follows easily from (6.7).
Corollary 3 Let the function $f(z)$ defined by (1.7) be in the class $Q_{m}^{n}(\Phi, \Psi, \lambda, A, B)$, then $f(z)$ is convex of order $\rho(0 \leq \rho<1)$ in $|z|<r_{3} \equiv r_{3}(m, n, \Phi, \Psi, \lambda, A, B, \rho)$, where

$$
r_{3}=\inf _{k}\left[\frac{(1-\rho)\left[(1-\lambda) \sigma(m, k) \beta_{k}+\lambda \sigma(n, k) \gamma_{k}\right](1+B)}{k(k-\rho)(B-A)}\right]^{1 / k-1}, \quad(k \geq 2)
$$

The result is sharp with the extremal function $f(z)$ given by (2.4).

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# New results in the stability study of non-autonomous evolution equations in Banach spaces ${ }^{1}$ 

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#### Abstract

The work extends some previous results of the first author, results concerning the Lyapunov stability of the zero solution of some nonautonomous nonlinear evolution equation.


AMS Subject Classification: 35B35, 35K90
Key Words and Phrases: abstract evolution equations, stability

## 1 Introduction

In previous works of the first author [5], [6], a nonlinear evolution problem in a Banach space $(X,\| \|)$, of the form

$$
\begin{gather*}
\dot{x}=A x+R(t, x),  \tag{1}\\
x(0)=a,
\end{gather*}
$$

is considered, and the existence of the solution in a neighborhood of $x=0$ as well as the Lyapunov stability of the null solution are studied via the implicit operator theorem.

In [5] the operator $A: D(A) \subset X \rightarrow X(D(A)$ dense in $X)$ is a closed linear operator that generates a strongly continuous, exponentially decreasing semigroup $T(t)$ on $X$, while the nonlinear operator $R$ is continuous, $R(t, 0)=0$ for all $t \in \mathbb{R}^{+}$ and for some $\beta>0, C>0$ the inequality
$\left\|R\left(t, x_{1}\right)-R\left(t, x_{2}\right)\right\| \leq C \max ^{\beta}\left(\left\|x_{1}\right\|,\left\|x_{2}\right\|\right)\left\|x_{1}-x_{2}\right\|$ holds for all $t \in \mathbb{R}^{+}$, and $x_{1}, x_{2}$ in a centered in 0 ball of $X$.

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The hypothesis on the nonlinear operator $R$ is weakened in [6], where the constant $C$ in the above inequality is replaced with a continuous function $C($.$) with exponential$ growth. In both papers above it is proved that, for an initial condition with small enough norm, the problem has an unique mild solution (solution of some integral equations associated to the differential abstract equation) and the null solution is exponentially stable in the class of mild solutions. If, moreover, a Hölder condition in $t$ is imposed to $R$, then the above conclusions hold for the classical solutions also.

## 2 Problem 1.

We consider the nonautonomous problem

$$
\begin{equation*}
\dot{x}=A x+R(t, x), \quad t \in \mathbb{R}, \tag{2}
\end{equation*}
$$

with the initial condition

$$
\begin{equation*}
x(s)=a \tag{3}
\end{equation*}
$$

for some fixed $s \in \mathbb{R}$.
In [7] the corresponding autonomous problem $(R(t, x)=R(x))$ was considered.

### 2.1 Hypotheses

I. $A: D(A) \subset X \mapsto X$, is a closed linear operator, $D(A)$ is dense in $X$.
II. The null space $\mathcal{V}=N(A)$ of the operator $A$ is nontrivial, closed, and $X$ may be written as the direct sum of $\mathcal{V}$ and $\mathcal{U}=R(A)(X=N(A) \oplus R(A))$.

Hence for any $x \in X$, there is an unique decomposition $x=u+v$ with $u \in \mathcal{U}, v \in$ $\mathcal{V}$. The mapping $P$ from $X$ onto $\mathcal{U}$, given by $P(x)=u$, is continuous.

The restriction of the operator $A$ to the subspace $\mathcal{U}$ is the generator of the exponentially decreasing semigroup $U(t)$ of class $C_{0}$ (there exist the constants $M>0$ and $\alpha>0$ such that for all $t \in \mathbb{R}^{+}=[0,+\infty)$ the inequality $\|U(t)\| \leq M \exp (-\alpha t)$ is fulfilled).
III. The nonlinear mapping $R$, defined on the Cartesian product of $\mathbb{R}$ with a neighborhood of 0 in $X$, is continuous, $R(t, 0)=0$ for all $t \in \mathbb{R}$ and there is a $\beta>0$ and a continuous function $C(t)>0$ such that for all $t \in \mathbb{R}$ and $x_{1}, x_{2}$ in a neighborhood of 0 in $X$, the inequality

$$
\begin{equation*}
\left\|R\left(t, x_{1}\right)-R\left(t, x_{2}\right)\right\| \leq C(t) \max \left(\left\|x_{1}\right\|^{\beta},\left\|x_{2}\right\|^{\beta}\right)\left\|x_{1}-x_{2}\right\| \tag{4}
\end{equation*}
$$

holds, where $C(\cdot) \in L^{1}(\mathbb{R})$ hence there is a $\widetilde{C}>0$ such that, for every $s \in \mathbb{R}$

$$
\int_{s}^{\infty}|C(\theta)| d \theta \leq \widetilde{C}
$$

We remark that the kernel of $A$ may be infinite dimensional, hence $A$ is not necessarily a Fredholm operator. We do not insist here on the conditions on $A$ such
that the above hypotheses hold. A comprehensive study of linear semigroups is given in [1].

By hypothesis II, $I-P$ maps $X$ on $N(A)$ and it is continuous. We denote, for simplicity, $\|P\|_{\mathcal{L}(X)}$ by $\|P\|$.

We set in (2) $x(t)=u(t)+v(t)$, where $u(t)=P x(t), v(t)=(I-P) x(t)$ and we project the problem (2), (3) on $U$ and $V$. We get the system of Cauchy problems

$$
\begin{gather*}
\dot{u}=P A u+P R(t, u+v), \quad u(s)=P a  \tag{5}\\
\dot{v}=(I-P) A u+(I-P) R(t, u+v), \quad v(s)=(I-P) a . \tag{6}
\end{gather*}
$$

We consider the system of integral equations

$$
\begin{gather*}
u(t)=U(t-s) P a+\int_{s}^{t} U(t-\theta) P R(\theta, u(\theta)+v(\theta)) d \theta  \tag{7}\\
v(t)=(I-P) a+\int_{s}^{t}(I-P)[A u(\theta)+R(\theta, u(\theta)+v(\theta))] d \theta . \tag{8}
\end{gather*}
$$

If the pair $(u(t), v(t))$ is a solution of the integral equations (7)-(8) then $x(t)=$ $u(t)+v(t)$ is called the mild solution of the Cauchy problem (2), (3).

Hypothesis II implies $(I-P) A=0$ and the equations above become

$$
\begin{gather*}
u(t)=U(t-s) P a+\int_{s}^{t} U(t-\theta) P R(\theta, u(\theta)+v(\theta)) d \theta  \tag{9}\\
v(t)=(I-P) a+\int_{s}^{t}(I-P) R(\theta, u(\theta)+v(\theta)) d \theta \tag{10}
\end{gather*}
$$

We set $t=\tau+s$ and we perform the change of variables $\theta^{\prime}=\theta-s$, to obtain

$$
\begin{gather*}
u(\tau+s)=U(\tau) P a+\int_{0}^{\tau} U\left(\tau-\theta^{\prime}\right) P R\left(\theta^{\prime}+s, u\left(\theta^{\prime}+s\right)+v\left(\theta^{\prime}+s\right)\right) d \theta^{\prime}  \tag{11}\\
v(\tau+s)=(I-P) a+\int_{0}^{\tau}(I-P) R\left(\theta^{\prime}+s, u\left(\theta^{\prime}+s\right)+v\left(\theta^{\prime}+s\right)\right) d \theta^{\prime} \tag{12}
\end{gather*}
$$

For every $s \in \mathbb{R}$, and any function $x:[s, \infty) \mapsto X$ we define the function $\left.x\right|_{s}$ : $[0, \infty) \mapsto X$, by $\left.x\right|_{s}(\theta)=x(s+\theta)$.

For a space $X_{1} \subseteq X$, we consider the space $\mathcal{C}_{b}\left([0, \infty), X_{1}\right)$ of continuous bounded functions defined on $[0, \infty)$ with values in $X_{1}$, with the supremum norm $\left(\|x\|_{0}=\right.$ $\left.\sup _{t \geq 0}\|x(t)\|\right)$, and denote it shortly by $\mathcal{C}_{b}\left(X_{1}\right)$. $t \geq 0$

We define the operators

$$
\begin{gathered}
D_{1}: X \mapsto C_{b}(\mathcal{U}), a \mapsto D_{1}(a), D_{1}(a)(\tau)=U(\tau) P a, \tau \geq 0 \\
D_{2}: X \mapsto C_{b}(\mathcal{V}), a \mapsto D_{2}(a), D_{2}(a)(\tau)=(I-P) a
\end{gathered}
$$

$$
F_{1 s}: \mathcal{C}_{b}(\mathcal{X}) \mapsto \mathcal{C}_{b}(\mathcal{U}), F_{1 s}(x)(\tau)=\int_{0}^{\tau} U(\tau-\theta) P R(s+\theta, x(\theta)) d \theta
$$

where $x=u+v$, and

$$
F_{2 s}: \mathcal{C}_{b}(\mathcal{X}) \mapsto \mathcal{C}_{b}(\mathcal{V}), \quad F_{2 s}(x)(\tau)=\int_{0}^{\tau}(I-P) R(s+\theta, x(\theta)) d \theta
$$

For these operators the following relations hold:

$$
\begin{gathered}
\left\|D_{1}\right\|_{\mathcal{L}\left(X, \mathcal{C}_{b}(\mathcal{U})\right)} \leq M\|P\| \\
\left\|D_{2}\right\|_{\mathcal{L}\left(X, \mathcal{C}_{b}(\mathcal{V})\right)} \leq 1+\|P\| \\
\left\|F_{1 s}(x)(\tau)\right\| \leq \int_{0}^{\tau} M e^{-\alpha(\tau-\theta)}\|P\| C(s+\theta)\|x(\theta)\|^{\beta+1} d \theta \leq M \widetilde{C}\|P\|\|x\|_{0}^{\beta+1} \\
\left\|F_{2 s}(x)(\tau)\right\| \leq \widetilde{C}(1+\|P\|)\|x\|_{0}^{\beta+1}
\end{gathered}
$$

The two integral equations (11), (12) may be written as the equation

$$
\begin{equation*}
\left.x\right|_{s}=D(a)+F_{s}\left(\left.x\right|_{s}\right) \tag{13}
\end{equation*}
$$

where

$$
D=D_{1}+D_{2}, F_{s}=F_{1 s}+F_{2 s}
$$

Equation (13) may be regarded as a fixed point problem in $\mathcal{C}_{b}(X)$,

$$
\begin{equation*}
\phi=\Phi(a, \phi) \tag{14}
\end{equation*}
$$

where

$$
\begin{equation*}
\Phi(a, \phi)=D(a)+F_{s}(\phi) . \tag{15}
\end{equation*}
$$

We remark that $\phi^{*}(a)$ is a fixed point of $\Phi(a,$.$) if and only if the function x(\cdot ; s, a)$ : $[s, \infty) \mapsto X$, defined by $x(s+\theta ; s, a)=\phi^{*}(\theta, a), \theta \geq 0$, is a mild solution of problem (2), (3).

### 2.2 Results

Theorem 1. There are two positive numbers $r_{0}, r_{1}$, such that for $\|a\| \leq r_{0}$, the mapping $\Phi(a, \cdot)$ given by (15) is an uniform (with respect to a) contraction from $\bar{B}\left(0, r_{1}\right) \subset \mathcal{C}_{b}(X)$ to itself.

Proof. We have, for any $\phi_{1}, \phi_{2} \in \mathcal{C}_{b}(X)$

$$
\begin{aligned}
& \left\|\Phi\left(a, \phi_{1}\right)-\Phi\left(a, \phi_{2}\right)\right\|_{0}=\left\|F_{s}\left(\phi_{1}\right)-F_{s}\left(\phi_{2}\right)\right\|_{0} \leq \\
& \leq\left\|F_{1 s}\left(\phi_{1}\right)-F_{1 s}\left(\phi_{2}\right)\right\|_{0}+\left\|F_{2 s}\left(\phi_{1}\right)-F_{2 s}\left(\phi_{2}\right)\right\|_{0}
\end{aligned}
$$

We have

$$
\left\|F_{1 s}\left(\phi_{1}\right)-F_{1 s}\left(\phi_{2}\right)\right\|_{0} \leq \int_{s}^{t} M e^{-\alpha(t-\theta)}\left\|P R\left(s+\theta, \phi_{1}(\theta)\right)-P R\left(s+\theta, \phi_{2}(\theta)\right)\right\| d \theta \leq
$$

$$
\begin{gather*}
\leq \int_{s}^{t} M e^{-\alpha(t-\theta)}\|P\| C(s+\theta) \max \left(\left\|\phi_{1}(\theta)\right\|^{\beta},\left\|\phi_{2}(\theta)\right\|^{\beta}\right)\left\|\phi_{1}(\theta)-\phi_{2}(\theta)\right\| d \theta \leq \\
\leq M \widetilde{C}\|P\| \max \left\{\left\|\phi_{1}\right\|_{0}^{\beta},\left\|\phi_{2}\right\|_{0}^{\beta}\right\}\left\|\phi_{1}-\phi_{2}\right\|_{0},(\forall) \phi_{1}, \phi_{2} \in \mathcal{C}_{b}(X) \tag{16}
\end{gather*}
$$

and

$$
\begin{gathered}
\left\|F_{2 s}\left(\phi_{1}\right)-F_{2 s}\left(\phi_{2}\right)\right\|_{0} \leq \int_{0}^{\tau}\left\|(I-P) R\left(s+\theta, \phi_{1}(\theta)\right)-(I-P) R\left(s+\theta, \phi_{2}(\theta)\right)\right\| d \theta \\
\leq \int_{0}^{\tau}(1+\|P\|)|C(s+\theta)| d \theta \max \left\{\left\|\phi_{1}\right\|_{0}^{\beta},\left\|\phi_{2}\right\|_{0}^{\beta}\right\}\left\|\phi_{1}-\phi_{2}\right\|_{0} \\
\leq \widetilde{C}(1+\|P\|) \max \left\{\left\|\phi_{1}\right\|_{0}^{\beta},\left\|\phi_{2}\right\|_{0}^{\beta}\right\}\left\|\phi_{1}-\phi_{2}\right\|_{0} .
\end{gathered}
$$

For $\phi_{1}, \phi_{2} \in B(0, r) \subset \mathcal{C}_{b}(X)$, we have

$$
\left\|F_{s}\left(\phi_{1}\right)-F_{s}\left(\phi_{2}\right)\right\|_{0} \leq \widetilde{C}[(M+1)\|P\|+1] r^{\beta}\left\|\phi_{1}-\phi_{2}\right\|_{0}
$$

Let $r_{1}>0$ be such that $r_{1}<\{\widetilde{C}[(M+1)\|P\|+1]\}^{-1 / \beta}$. On the sphere $\bar{B}\left(0, r_{1}\right) \subset$ $\mathcal{C}_{b}(X)$, the mapping $\Phi(a, \cdot)$ is an uniform (with respect to $a$ ) contraction.

If $\|a\| \leq r_{0},\|\phi\|_{0} \leq r_{1}$, then

$$
\|\Phi(a, \phi)\|_{0} \leq[(M+1)\|P\|+1] r_{0}+\widetilde{C}[(M+1)\|P\|+1] r_{1}^{\beta+1}
$$

By imposing to this last quantity to be less than $r_{1}$, we find the restriction

$$
\|a\| \leq r_{0}:=\frac{r_{1}}{(M+1)\|P\|+1}\left\{1-\widetilde{C}[(M+1)\|P\|+1] r_{1}^{\beta}\right\}
$$

Hence, for $\|a\| \leq r_{0}, \Phi$ is an uniform contraction on $\bar{B}\left(0, r_{1}\right)$.
The uniform Banach contraction principle implies that, for $\|a\| \leq r_{0}$, the fixed point problem (14) has an unique solution $\varphi^{*}(a) \in \mathcal{C}_{b}(X)$.

Hence the function $x(\cdot ; s, a)$ given by $x(s+\theta ; s, a)=\varphi^{*}(\theta, a), \theta \geq 0$ is a mild solution of problem (2), (3). Since $\Phi$ is an uniform contraction with respect to $a$, the fixed point $\varphi^{*}$ is continuous with respect to $a$. Hence the function that maps $a \in X$ to $x(s+\cdot ; s, a) \in \mathcal{C}_{b}(X)$ is continuous. From here and from $x(. ; s, 0)=0$, the Lyapunov stability of the null solution in the class of mild solutions follows, where the Lyapunov stability is understood in the following sense

Definition 1. [2] A classical (resp. -mild) solution $x(. ; s, a)$ of problem (2), (3) is called stable if for every $\epsilon>0$ and every $s^{\prime}>s$ there is a $\delta=\delta\left(\epsilon, s^{\prime}\right)$ such that for every $y \in X$ with $\left\|y-x\left(s^{\prime} ; s, a\right)\right\| \leq \delta$, the classical (resp. -mild) solution $x\left(\cdot ; s^{\prime}, y\right)$ exists, is defined on $\left[s^{\prime}, \infty\right)$, and

$$
\left\|x(t ; s, a)-x\left(t ; s^{\prime}, y\right)\right\|<\epsilon
$$

for every $t \geq s^{\prime}$.
If a Hölder condition with respect to time is added on $R$, then we obtain the existence of classical solutions and also the stability of the 0 solution in the class of classical solutions.

Remark. If we assume also that $R(t, v)=0$ for every $v \in \mathcal{V}$, then by observing that in this case $\|R(\theta, u(\theta)+v(\theta))\| \leq C(\theta)\|u(\theta)\|^{\beta+1}$ holds, and by the reasoning of [5] we obtain that $\|u(t)\|$ tends exponentially to 0 when $t \rightarrow \infty$.

## 3 Problem 2.

We consider a perturbation of Problem 1, that is

$$
\begin{equation*}
\dot{x}=A x+B x+R(t, x), \quad t \in \mathbb{R}, \tag{17}
\end{equation*}
$$

with the initial condition

$$
\begin{equation*}
x(s)=a \tag{18}
\end{equation*}
$$

Here $B$ is a linear bounded operator $B: X \mapsto R(A)$ with norm such that

$$
\begin{equation*}
L:=\frac{M}{\alpha}\|P\|\|B\|_{\mathcal{L}(X)}<1 \tag{19}
\end{equation*}
$$

Due to the range of $B$, only the first equation of the projected equations differs from those of the previous section. The projected equations are, in this case,

$$
\begin{gather*}
\dot{u}=P A u+P B(u+v)+P R(t, u+v), \quad u(s)=P a,  \tag{20}\\
\dot{v}=(I-P) R(t, u+v), \quad v(s)=(I-P) a . \tag{21}
\end{gather*}
$$

We consider the integral equations

$$
\begin{gather*}
u(t)=U(t-s) P a+\int_{s}^{t} U(t-\theta) P B u(\theta) d \theta+\int_{s}^{t} U(t-\theta) P R(\theta, u(\theta)+v(\theta)) d \theta,  \tag{22}\\
v(t)=(I-P) a+\int_{s}^{t}(I-P) R(\theta, u(\theta)+v(\theta)) d \theta \tag{23}
\end{gather*}
$$

We say that a function $x(\cdot)=u(\cdot)+v(\cdot)$ is an $A$-mild solution of equation (17) if $(u(\cdot), v(\cdot))$ is the solution of $(22)-(23)$. We named this solution $A$-mild because the semigroup generated only by $A$ (and not by $A+B$ ) is used in the first integral equation. We proceed as for Problem 1, only that we define also the operator

$$
E: \mathcal{C}_{b}(\mathcal{X}) \mapsto \mathcal{C}_{b}(\mathcal{U}), E(x)(t)=\int_{s}^{t} U(t-\theta) P B(x(\theta)) d \theta
$$

for which

$$
\|E(\phi)\|_{0} \leq L\|\phi\|_{0}
$$

with $L$ given by (19) and

$$
\left\|E\left(\phi_{1}\right)-E\left(\phi_{2}\right)\right\|_{0} \leq L\left\|\phi_{1}-\phi_{2}\right\|_{0} .
$$

The integral equations above are equivalent with the fixed point problem in $\mathcal{C}_{b}(X)$

$$
\phi=\Phi(a, \phi)
$$

where

$$
\begin{equation*}
\Phi(a, \phi)=D(a)+E(\phi)+F_{s}(\phi) . \tag{24}
\end{equation*}
$$

Theorem 2. There are two positive numbers $r_{0}, r_{1}$, such that for $\|a\| \leq r_{0}$, the mapping $\Phi(a, \cdot)$ given by (24) is an uniform (with respect to a) contraction from $\bar{B}\left(0, r_{1}\right) \subset \mathcal{C}_{b}(X)$ to itself.

Proof. The computations in the proof of Theorem 1 and the properties of $E$ lead to

$$
\left\|\Phi\left(a, \phi_{1}\right)-\Phi\left(a, \phi_{2}\right)\right\|_{0} \leq\left\{\widetilde{C}[(M+1)\|P\|+1] r^{\beta}+L\right\}\left\|\phi_{1}-\phi_{2}\right\|_{0}
$$

for $\left\|\phi_{i}\right\|_{0} \leq r, i=1,2$.
We choose a $r_{1}>0$ such that $r_{1}<\left(\frac{1-L}{\widetilde{C}[(M+1)\|P\|+1]}\right)^{1 / \beta}$.
On the closed sphere $\|\phi\|_{0} \leq r_{1}$ the mapping $\Phi$ is a contraction (uniform with respect to $a$ ). We have

$$
\left\|\Phi\left(a, \phi_{1}\right)\right\|_{0} \leq[(M+1)\|P\|+1]\|a\|+\widetilde{C}[(M+1)\|P\|+1] r_{1}^{\beta+1}+L r_{1}
$$

We impose the condition

$$
[(M+1)\|P\|+1]\|a\|+\widetilde{C}[(M+1)\|P\|+1] r_{1}^{\beta+1}+L r_{1} \leq r_{1}
$$

and we find that this condition is satisfied for

$$
\|a\| \leq r_{0}:=\frac{1}{(M+1)\|P\|+1}\left\{1-L-\widetilde{C}[(M+1)\|P\|+1] r_{1}^{\beta}\right\} r_{1} .
$$

The assertion of the theorem follows.
From this point, with the same reasonings as in the preceding section, we obtain that, for $\|a\| \leq r_{0}, r_{0}$ defined in the proof of Theorem 2, the problem (17), (18) has an unique $A$-mild solution $x(\cdot ; s, a)$ in the space of continuous bounded functions defined from $[s, \infty)$ to $X$. Moreover, from the continuity of the function that maps $a \in X$ to $x(s+\cdot ; s, a) \in \mathcal{C}_{b}(X)$ the Lyapunov stability of the null solution in the class of $A$-mild solutions follows.

If the nonlinear map $R(.,$.$) satisfies a Hölder condition with respect to time, then$ the above assertions are valid in the classical sense also.

Remark. If we assume also that $R(t, v)=0$ for $v \in \mathcal{V}$, then by observing that $\|R(\theta, u(\theta)+v(\theta))\| \leq C(\theta)\|u(\theta)\|^{\beta+1}$, and by reasonings similar to those of Section 4 of [3] we obtain that $\|u(t)\| \rightarrow 0$ when $t \rightarrow \infty$.

## 4 Problem 3.

We consider again equation (2) with condition (3), but here we assume that only Hypothesis I on $A$ and Hypothesis III on $R$ are fulfilled (with no special hypotheses on the kernel of $A$ ). Moreover, we assume that $A$ generates a bounded semigroup $\{T(t)\}_{t \geq 0}$ of operators on $X$, that is for a $M>0,\|T(t)\|_{\mathcal{L}(X)} \leq M, \forall t \geq 0$.

We consider the integral equation

$$
\begin{equation*}
x(t)=T(t-s) a+\int_{s}^{t} T(t-\theta) R(\theta, x(\theta)) d \theta \tag{25}
\end{equation*}
$$

A solution of this equation is called a mild solution of problem (2), (3). Obviously, any classical solution is a mild solution of this problem. As before, we successively transform the integral equation (25) into a fixed point problem:

$$
x(s+\tau)=T(\tau) a+\int_{0}^{\tau} T\left(\tau-\theta^{\prime}\right) R\left(s+\theta^{\prime}, x\left(s+\theta^{\prime}\right)\right) d \theta^{\prime}
$$

hence

$$
\begin{equation*}
\left.x\right|_{s}(\tau)=T(\tau) a+\int_{0}^{\tau} T\left(\tau-\theta^{\prime}\right) R\left(s+\theta^{\prime},\left.x\right|_{s}\left(\theta^{\prime}\right)\right) d \theta^{\prime} \tag{26}
\end{equation*}
$$

Now, we define $D: X \mapsto \mathcal{C}_{b}(X), a \rightarrow D(a)(),. D(a)(\tau)=T(\tau) a$, and $F_{s}: \mathcal{C}_{b}(X) \mapsto$ $\mathcal{C}_{b}(X)$

$$
F_{s}(\phi)=\int_{0}^{\tau} T\left(\tau-\theta^{\prime}\right) R\left(s+\theta^{\prime}, \phi\left(\theta^{\prime}\right)\right) d \theta^{\prime}
$$

for which we have $\|D(a)\|_{0} \leq M\|a\|$, and, respectively

$$
\left\|F_{s}(\phi)\right\|_{0} \leq \int_{0}^{\tau} M C\left(s+\theta^{\prime}\right)\left\|\phi\left(\theta^{\prime}\right)\right\|^{\beta+1} d \theta^{\prime} \leq M \widetilde{C}\|\phi\|_{0}^{\beta+1}
$$

hence $F_{s}$ takes indeed values in $\mathcal{C}_{b}(X)$. Moreover,

$$
\begin{gathered}
\left\|F_{s}\left(\phi_{1}\right)-F_{s}\left(\phi_{2}\right)\right\|_{0} \leq \int_{0}^{\tau} M C\left(s+\theta^{\prime}\right) \max \left(\left\|\phi_{1}\left(\theta^{\prime}\right)\right\|^{\beta},\left\|\phi_{2}\left(\theta^{\prime}\right)\right\|^{\beta}\right)\left\|\phi_{1}\left(\theta^{\prime}\right)-\phi_{2}\left(\theta^{\prime}\right)\right\| d \theta^{\prime} \leq \\
\leq M \widetilde{C} \max \left(\left\|\phi_{1}\right\|_{0}^{\beta},\left\|\phi_{2}\right\|_{0}^{\beta}\right)\left\|\phi_{1}-\phi_{2}\right\|_{0}
\end{gathered}
$$

We define the mapping $\Phi: X \times \mathcal{C}_{b}(X) \mapsto \mathcal{C}_{b}(X)$, given by

$$
\begin{equation*}
\Phi(a, \phi)=D(a)+F_{s}(\phi) \tag{27}
\end{equation*}
$$

and remark that equation (26) is equivalent to the fixed point problem in $\mathcal{C}_{b}(X)$

$$
\begin{equation*}
\phi=\Phi(a, \phi) \tag{28}
\end{equation*}
$$

Now consider a positive number $r_{1}$ such that $M \widetilde{C} r_{1}^{\beta}<1$. On the ball $B\left(0, r_{1}\right)$ the mapping $\Phi$ is a contraction (uniform with respect to $a$ ). Now we take an $a$ such that

$$
M\|a\|+M \widetilde{C} r_{1}^{\beta+1} \leq r_{1} \Leftrightarrow\|a\| \leq r_{0}:=r_{1}\left(1-M \widetilde{C} r_{1}^{\beta}\right) / M
$$

We thus proved
Theorem 3. There are two positive numbers $r_{0}, r_{1}$, such that for $\|a\| \leq r_{0}$, the mapping $\Phi(a, \cdot)$ given by (27) is an uniform (with respect to a) contraction from $\bar{B}\left(0, r_{1}\right) \subset \mathcal{C}_{b}(X)$ to itself.

Hence, if $\|a\|<r_{0}$, there is an unique solution of equation (28) in $\bar{B}\left(0, r_{1}\right)$. Thus an unique mild solution $x(. ; s, a)$ of problem (2), (3) exists. Then, the local Lyapunov stability of the mild 0 solution follows from the continuity of the function $a \in X \rightarrow$ $x(s+\cdot ; s, a) \in \mathcal{C}_{b}(X)$.

If we add a Hölder condition in $t$ on $R$ the local stability of 0 as a classical solution follows also.

### 4.1 Problem 4.

We show in this section that similar results may be obtained also in the case when the linear operator is a function of time, that is for the equation

$$
\begin{equation*}
\dot{x}=A(t) x+R(t, x), t \in \mathbb{R} \tag{29}
\end{equation*}
$$

with the initial condition

$$
x(s)=a
$$

Here $A(t): D \subset X \mapsto X$ are closed linear operators, defined on the domain $D$ (independent of $t$ ), dense in $X$, and we assume that the problem

$$
\begin{equation*}
\dot{u}=A(t) u, u(s)=a \tag{30}
\end{equation*}
$$

has an unique classical solution, that is a function $u \in \mathcal{C}([s, \infty), X) \cap \mathcal{C}^{1}([s, \infty), X)$ such that $u(t) \in D$ for $t>s$ and (30) is satisfied. In this situation by the relation

$$
U(t, s) u(s)=u(t)
$$

an evolution family is generated.
An evolution family is a family of linear bounded operators $\{U(t, s), t \geq s, s \in \mathbb{R}\}$ that satisfy

$$
U(t, r) U(r, s)=U(t, s), U(s, s)=I
$$

We say that an evolution family is strongly continuous if the mapping $(t, s) \rightarrow U(t, s)$ is strongly continuous on the set $t, s \in \mathbb{R}, t>s$.

The most difficult problem concerning this case is that of giving sufficient conditions in order that the linear problem attached to the above problem generates such an evolution family. This is not the subject of this paper and we indicate the work [4] and the references therein.

The evolution family is named exponentially bounded if the inequality

$$
\|U(t, s)\|_{\mathcal{L}(X)} \leq M e^{\omega(t-s)}
$$

$M \geq 1, \omega \in \mathbb{R}$, holds. If $\omega<0$, the evolution family is named exponentially stable.
We first assume that the evolution family generated by (30) is exponentially stable hence there is a positive $\alpha$ such that

$$
\begin{equation*}
\|U(t, s)\|_{\mathcal{L}(X)} \leq M e^{-\alpha(t-s)} \tag{31}
\end{equation*}
$$

In order to study the existence of the solution of problem (29), we consider the integral equation

$$
\begin{equation*}
x(t)=U(t, s) a+\int_{s}^{t} U(t, \theta) R(\theta, x(\theta)) d \theta \tag{32}
\end{equation*}
$$

Any classical solution of problem (29) is a solution of (32). A solution of (32) is called a mild solution of (29). From this point, formally, the reasoning may be lead as in [6], since only the norm estimates and the hypotheses on $R$ will be used.

In [6] the hypotheses on $R$ differ from Hypothesis III from the beginning of this work only in what concerns the function $C(\cdot)$. More precisely, it is assumed that one of the inequalities

$$
\begin{gathered}
C(t) e^{-\gamma \beta} \leq C^{*}, \text { for a } \gamma \in(0, \alpha), C^{*}>0 \\
\int_{0}^{\infty} C(s) e^{-\alpha \beta s} d s<\infty
\end{gathered}
$$

holds (hence $C($.$) may be unbounded).$
As in [6], the problem may be brought to the form of a fixed point problem in the space $\mathcal{C}_{\gamma}(X)=\left\{x:[0, \infty) \mapsto X \mid \sup _{[0, \infty)}\|x(t)\| e^{\gamma t}<\infty\right\}$. It follows that for sufficient small $\|a\|$ there is an unique mild solution of (29) and that the null solution is locally asymptotically stable (in the class of mild solutions), in the sense given by

Definition 2. [2] A classical (resp. -mild) solution $x(. ; s, a)$ of problem (2), (3) is called asymptotically stable if it is stable and, for every $s^{\prime}>s$, there is a $\delta=\delta\left(s^{\prime}\right)>0$ such that for $y \in X$ with $\left\|y-x\left(s^{\prime} ; s, a\right)\right\|<\delta$, the classical (resp. -mild) solution $x\left(\cdot ; s^{\prime}, y\right)$ exists, is defined on $\left[s^{\prime}, \infty\right)$, and

$$
\left\|x(t ; s, a)-x\left(t ; s^{\prime}, y\right)\right\| \rightarrow 0
$$

as $t \rightarrow \infty$.
If we assume that the evolution family is only bounded $(\alpha=0)$, then with the same reasonings as in Section 3 of the present paper, we obtain that for sufficient small $\|a\|$ there is an unique mild solution of (29) and that the null solution is Lyapunov stable (in the class of mild solutions).

If the function $R(\cdot, \cdot)$ satisfies also a Hölder condition in $t$ the above conclusions are valid in the classical sense also.

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