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Functional Calculus

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Dr. Hammad Khalil is an assistant professor of mathematics in the Department of Mathematics at the University of Education, Lahore. He completed his Ph.D. in applied mathematics under the supervision of Professor Dr. Rahmat Ali Khan. He has published numerous articles in various international journals and serves on a number of editorial boards. His research areas include functional analysis, approximation theory, chaos theory and nonlinear dynamics.

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## Preface

Calculus is the elementary subject of applied analysis and its study includes a rich variety of functions and their behavior. This book brings together a range of different concepts from across the wide spectrum of the concept of calculus.

The book is in two sections. The first deals with advances in analysis and the second with the application of some results of functional calculus to applied problems.

The first section opens with an analysis of logarithmic potential transform. The singular values of this transform are discussed on the Poincaré disk. This potential can be used to illustrate some of the important features of field theory such as dimensional regularization and renormalization. Although most recent textbooks do not discuss this potential in detail, the calculations to demonstrate some of its unique features are quite simple. The bound state energy of this logarithmic potential is obtained through the uncertainty principle, phase space quantization and the Hellmann-Feynman theorem.

In the second chapter of the first section, the authors define and prove new Tauberian theorems under triple statistically Norlund-Cesaro summability. Some theorems, lemmas and corollaries can be defined and proved similarly by using the ( $1,0,0$ ), $(0,1,0)$ and $(0,0,1)$ summability method. Although Tauberian theorems for single sequences of a single variable are well established, they remain in their infancy for triple sequences.

The final chapter of the first section is devoted to the study of the Calderon operator, which is the sum of the Hardy averaging operator and its adjoint and plays an important role in the theory of real interpolation. On the other hand, the Hilbert operator arises from the continuous version of Hilbert's inequality. Both operators appear in different contexts and have numerous applications within the harmonic analysis. In this chapter, the authors briefly review the Calderon and Hilbert operators, showing some of the most relevant results within the functional analysis and presenting recent results on these operators within Fourier analysis.

The second section of the book collects some results from applied analysis The first chapter deals with the study of heat transfer development of titanium oxide nanofluid of platelet-shaped nanoparticles over a vertical stretching cylinder. A set of nonlinear equations is obtained using suitable transformation on the governing equations which are then solved with numerical scheme BVP4C. The results obtained are interpreted graphically and numerically. The effects of Prandtl, Eckert and unsteadiness parameters on temperature distribution are depicted, and skin friction and Nusselt number are also computed. In the second chapter, a nonlinear response of the follower motion is simulated at different cam speeds, different coefficients of restitution and different internal distances of the follower guide from inside. The nonlinear response of the follower is employed to investigate the chaotic phenomenon in the cam follower system in
the presence of follower offset. The numerical results are achieved using Solid Works software. The chaos phenomenon is detected using Poincaré maps with phase-plane portraits, the largest Lyapunov exponent parameter, and a bifurcation diagram. The largest Lyapunov exponent has its maximum value when the follower offsets to the right, and its minimum value when the follower offsets to the left. The chaotic phenomenon in cam follower systems when the follower offsets to the left is greater than the chaotic phenomenon when the follower offsets to the right. The final chapter investigates the decision fusion problem for large-scale sensor networks associated with the Internet of Things and artificial intelligence. The sensor networks discussed are those with unavoidable transmission channel interference and non-ideal channels that are prone to errors. A generalized algorithm is proposed that enables decision fusion rules to be designed for large-scale sensor networks and can at the same time search for the optimal sensor rules and the optimal fusion rule. Finally, numerical examples show the effectiveness of the new algorithms for large-scale sensor networks with non-ideal channels.

Section 1
Advances in Theory

## Chapter 1

# The Singular Values of the Logarithmic Potential Transform on Bound States Spaces of Landau Hamiltonians on the Poincaré Disk 

M'hamed Elomari and Ali El Mfadel


#### Abstract

In the present manuscript, we prove that the singular numbers of the Cauchy transform $\mathcal{L}_{\sigma}[f](z)=-\frac{1}{\pi} \int_{\mathbb{D}} \frac{f(\xi)}{\xi-z} \log \left(\frac{1}{z-\xi \mid}\right)(1-\xi \bar{\xi})^{\sigma-2} d \mu(\xi)$ (2) defined on the space $L^{2, \sigma}(\mathbb{D})$ of complex-valued measurable functions, which are $(1-\xi \bar{\xi})^{\sigma-2} d \mu(\xi)$-square integrable on $\mathbb{D}$ where $\sigma>1$ is a fixed parameter, are asymptotically $\approx C \sqrt{k^{m-4 \nu+1}}$, as $k \rightarrow \infty$ where $C$ is a constant.


Keywords: the logarithmic potential transform, the singular values, Cauchy transform.

## 1. Introduction

Let $\mathbb{D}$ be the complex unit disk endowed with its Lebesgue measure $\mu$ and let $\partial \mathrm{ID}$ be its boundary. Denote by $L^{2}(\mathbb{D}, d \mu)$ the space of complex-valued measurable functions, which are $d \mu$ square integrable on $\mathbb{D}$. The logarithmic potential transform: $L^{2}(\mathbb{D}) \rightarrow L^{2}(\mathbb{D})$ is defined by

$$
\begin{equation*}
\mathcal{L}[f](z)=-\frac{1}{\pi} \int_{\mathbb{D}} \frac{f(\xi)}{\xi-z} \log \left(\frac{1}{|z-\xi|}\right) d \mu(\xi) . \tag{1}
\end{equation*}
$$

This operator is very important as the transformed Cauchy and it often appears in Analysis [1].

The dimensional analysis [1, 2] and scaling arguments form an integral part of theoretical physics to solve some important problems without doing much calculation.

The logarithmic potential in physics forms an interesting one as it provides some unusual predictions about the system. Moreover, this potential can be used suitably to illustrate some of the important features of field theory such as dimensional regularization and renormalization. In most of our textbooks, this potential is not discussed in detail; although the calculations are quite simple to demonstrate some of its unique features. We have obtained the bound state energy of this logarithmic potential
through uncertainty principle, phase space quantization, and the Hellmann-Feynman theorem.

In Ref. [3] the authors have been dealing with the restriction of $\mathcal{L}$ to the space $L_{a}^{2}(\mathbb{D})$ of analytic $\mu$-square integrable on $\mathbb{D}$. They precisely have considered the projection operator $P_{0}: L^{2}(\mathbb{D}) \rightarrow L_{a}^{2}(\mathbb{D})$ and they have proved that the singular values $\lambda_{k}$ of $\mathcal{L} P_{0}$, (which turn out to be eigenvalues of the operator $\sqrt{\left(\mathcal{L} \mathcal{P}_{0}\right)^{*}\left(\mathcal{L} P_{0}\right)}$ behave like $k^{-1}$ as $k$ goes to $\infty$. They also concluded that $\mathcal{L} P_{0}$ belongs to the Schatten class $S_{1, \infty}$.

Now, consider the following weighted logarithmic potential transform

$$
\begin{equation*}
\mathcal{L}_{\sigma}[f](z)=-\frac{1}{\pi} \int_{\mathbb{D}} \frac{f(\xi)}{\xi-z} \log \left(\frac{1}{|z-\xi|}\right)(1-\xi \bar{\xi})^{\sigma-2} d \mu(\xi) \tag{2}
\end{equation*}
$$

defined on the space $L^{2, \sigma}(\mathbb{D})$ of complex-valued measurable functions, which are $(1-\xi \bar{\xi})^{\sigma-2} d \mu(\xi)$-square integrable on $\mathbb{D}$ where $\sigma>1$ is a fixed parameter. We observe that the subspace $L_{a}^{2, \sigma}(\mathbb{D})$ of analytic functions on $\mathbb{D}$ and belonging to $L^{2, \sigma}(\mathbb{D})$ coincides with the eigenspace

$$
\begin{equation*}
\mathcal{A}_{0}^{\sigma}(\mathbb{D}):=\left\{\psi \in L^{2, \sigma}(\mathbb{D}), \Delta_{\sigma} \psi=0\right\}, \tag{3}
\end{equation*}
$$

of the second order differential operator

$$
\begin{equation*}
\Delta_{\sigma}:=-4(1-z \bar{z})\left((1-z \bar{z}) \frac{\partial^{2}}{\partial z \partial \bar{z}}-\sigma \bar{z} \frac{\partial}{\partial \bar{z}}\right), \tag{4}
\end{equation*}
$$

known as the $\sigma$-weight Maass Laplacian and its discrete eigenvalues are given by

$$
\begin{equation*}
\varepsilon_{m}:=4 m(\sigma-1-m), m=0,1,2, \ldots,\lfloor(\sigma-1) / 2\rfloor, \tag{5}
\end{equation*}
$$

with their corresponding eigenspaces

$$
\begin{equation*}
\mathcal{A}_{m}^{\sigma}(\mathbb{D}):=\left\{\psi \in L^{2, \sigma}(\mathbb{D}) \text { and } \Delta_{\sigma} \psi=\varepsilon_{m}^{\sigma} \psi\right\}, \tag{6}
\end{equation*}
$$

are here called generalized Bergman spaces since ...
After noticing that, we here deal with analogous questions as in Ref. [3] in the context of the weighted Cauchy transform (2) and for its restriction to the space $\mathcal{A}_{m}^{\sigma}(\mathbb{D})$. That is, we are concerned with the operator $\mathcal{C}_{\sigma} P_{m}^{\sigma}$ where $P_{m}^{\sigma}$ is the projection $L^{2, \sigma}(\mathbb{D}) \rightarrow \mathcal{A}_{m}^{\sigma}(\mathbb{D})$. The results achieved are as follows:

Firstly, we find that the singular values of $\mathcal{L}_{\sigma} P_{m}^{\sigma}$. For $k \neq m$, it can be expressed as

$$
\lambda_{k}=\sqrt{J_{1}+J_{2}+J_{3}},
$$

where

$$
\begin{gathered}
J_{1}=\left(\frac{(1+k-m)_{m}}{m!(k-m+1)}\right)^{2} \sum_{n=0}^{\infty} A_{n} \frac{\Gamma(2 n+2 k-2 m+6-1) \Gamma(4 \nu-2 m-1)}{\Gamma(2 n+2 k-4 m+4 \nu+6)}, \\
J_{2}=\left(\frac{\alpha_{k}^{\nu, m}}{2 \nu-m-1}\right)^{2} \sum_{n=0}^{\infty} A_{n} \frac{\Gamma(4 \nu-2 m-1) \Gamma(2 n+2)}{\Gamma(2 n+4 \nu-2 m+1)},
\end{gathered}
$$

and

$$
J_{3}=\frac{(1+k-m)_{m} \alpha_{k}^{\nu, m}}{m!(k-m+1)(2 \nu-m-1)}\left(\sum_{n=0}^{\infty} A_{n} \frac{\Gamma(k-m+2) \Gamma(4 \nu-2 m-1)}{\Gamma(4 \nu-k-3 m)}\right) .
$$

For $k=m$ can be expressed as

$$
\begin{equation*}
\lambda_{k}^{2}=\frac{\alpha_{k}^{\nu, m}(2(2 \nu-m)-1)}{8(\pi(2 \nu-m+1))} \sum_{n=0}^{\infty} \frac{B_{n}}{n+2 \nu-m}, \tag{7}
\end{equation*}
$$

where

$$
\begin{gathered}
B_{n}=\sum_{n=0}^{\infty} \frac{\Gamma(-m+1) \Gamma(2 \nu-m) \Gamma(2(\nu-m)+1)}{n!\Gamma(2(\nu-m) \Gamma(2 \nu-m+2)}, \\
\alpha_{k}^{\nu, m}=\frac{\Gamma(2) \Gamma(2(m-\nu)+1)}{\Gamma(m+1) \Gamma(2+m-2 \nu)} .
\end{gathered}
$$

Secondly, we show that these singular values behave like

$$
\lambda_{k} \sim C \sqrt{k^{m-4 \nu+1}}, \text { as } k \rightarrow \infty,
$$

where $C$ is a constant.
The paper is organized as follows: In Section 2, we review the definition of the weighted logarithmic potential transform, as well as some of its needed properties. Section 3 deals with some basic facts on the spectral theory of Mass Laplacians on the Poincaré disk. In Section 4, a precise description of the generalized Bergmann spaces is reviewed. Section 5 is devoted to the computation of the singular values of the weighted logarithmic potential transform. The asymptotic behavior of these singular values is established in Section 6.

## 2. The weighted logarithmic potential transform $\mathcal{L}_{\nu}$

### 2.1 The case $\nu=1$

Let $\mathbb{D}$ the complex unit disk endowed with its Lebesgue measure $\mu$ and let $\partial \mathrm{ID}$ its boundary denote by $L^{2}(\mathbb{D})$ the space of complex-valued measurable functions on $\mathbb{D}$ with finite norm

$$
\begin{equation*}
\|f\|=\int_{\mathbb{D}}|f(\xi)|^{2} d \mu(\xi) \tag{8}
\end{equation*}
$$

The Logarithmic Potential operator $\mathcal{L}: L^{2}(\mathbb{D}) \rightarrow L^{2}(\mathbb{D})$ is defined by

$$
\begin{equation*}
\mathcal{L}[f](z)=\int_{I D} f(\xi) \log \left(\frac{1}{|\xi-z|}\right) d \mu(\xi) . \tag{9}
\end{equation*}
$$

### 2.2 The case of $\nu \geq 1$

We fix a real parameter $\nu$ such that $2 \nu>1$ and we consider the following weighted logarithmic potential transform

$$
\begin{equation*}
\mathcal{L}_{\nu}[f](z)=\int_{\mathbb{D}} f(\xi) \log \left(\frac{1}{|\xi-z|}\right)(1-\xi \bar{\xi})^{2 \nu-2} d \mu(\xi) \tag{10}
\end{equation*}
$$

defined on the space $L^{2, \nu}(\mathbb{D})$ complex-valued measurable functions are $(1-\xi \bar{\xi})^{2 \nu-2} d \mu(\xi)$-square integrable on $\mathbb{D}$. As a convolution of $L^{2, \nu}$-functions with the compactly supported measure $\frac{(1-\xi)^{2 \nu-2}}{\xi} \mathbb{1}_{\mathrm{D}} d \mu(\xi) \mathcal{L}_{\nu}: L^{2, \nu}(\mathbb{D}) \rightarrow L^{2, \nu}(\mathbb{D})$ is obviously bounded. Moreover, it is not hard to show that $\mathcal{L}_{\nu}$ is in fact compact [4]. This raises a question concerning the spectral picture of $\mathcal{L}_{\nu}$.

## 3. The Landau Hamiltonian $H_{\nu}$ on the Poincaré disk $\mathbb{D}$

Let $\mathbb{D}=\{z \in \mathbb{C}, z \bar{z}<1\}$ be the complex unit disk with the Poincaré metric $d s^{2}=$ $4(1-z \bar{z})^{-2} d z d \bar{z} . \mathbb{D}$ is a complete Riemannian manifold with all sectional curvature equal -1 . It has an ideal boundary $\partial \mathrm{ID}$ identified with the circle $\{\omega \in \mathbb{C}, \omega \bar{\omega}=1\}$. One refers to points $\omega \in \partial \mathrm{ID}$ as points at infinity. The geodesic distance between two points $z$ and $w$ is given by

$$
\begin{equation*}
\cosh d(z, w)=1+\frac{2(z-w)(\bar{z}-\bar{w})}{(1-z \bar{z})(1-w \bar{w})} \tag{11}
\end{equation*}
$$

By Ref. [5] the Schrödinger operator on $\mathbb{D}$ with a constant magnetic field of strength proportional to $\nu>0$ can be written as:

$$
\begin{equation*}
\mathcal{L}_{\nu}:=-\left(1-|z|^{2}\right)^{2} \frac{\partial^{2}}{\partial z \partial \bar{z}}-\nu z\left(1-|z|^{2}\right) \frac{\partial}{\partial z}+\nu \bar{z}\left(1-|z|^{2}\right) \frac{\partial}{\partial \bar{z}}+\nu^{2}|z|^{2} . \tag{12}
\end{equation*}
$$

which is also called Maass Laplacian on the disk. A slight modification of $\mathcal{L}_{\nu}$ is given by the operator

$$
\begin{equation*}
H_{\nu}:=4 \mathcal{L}_{\nu}-4 \nu^{2} \tag{13}
\end{equation*}
$$

acting in the Hilbert space

$$
\begin{equation*}
L^{2,0}(\mathbb{D}):=\left\{\varphi: \mathbb{D} \rightarrow \mathbb{C}, \int_{\mathbb{D}}|\varphi(z)|^{2}\left(1-|z|^{2}\right)^{-2} d \mu(z)<+\infty\right\} \tag{14}
\end{equation*}
$$

For our purpose, we shall consider the unitary equivalent realization $\tilde{H}_{\nu}$ of the operator $H_{\nu}$ in the Hilbert space

$$
\begin{equation*}
L^{2, \nu}(\mathbb{D}):=\left\{\varphi: \mathbb{D} \rightarrow \mathbb{C}, \int_{\mathbb{D}}|\varphi(z)|^{2}\left(1-|z|^{2}\right)^{2 \nu-2} d \mu(z)<+\infty\right\}, \tag{15}
\end{equation*}
$$

which is defined by

$$
\begin{equation*}
\tilde{H}_{\nu}:=Q_{\nu}^{-1} H_{\nu} Q_{\nu}, \tag{16}
\end{equation*}
$$

where $Q_{\nu}: L^{2, \nu}(\mathbb{D}) \rightarrow L^{2,0}(\mathbb{D})$ is the unitary transformation defined by the map $\varphi \mapsto Q_{\nu}[\varphi]:=\left(1-|z|^{2}\right)^{-\nu} \varphi$. Different aspects of the spectral analysis of the operator $\tilde{H}_{\nu}$ have been studied by many authors. For instance, note that $\tilde{H}_{\nu}$ is an elliptic densely defined operator on the Hilbert space $L^{2, \nu}(\mathbb{D})$ and admits a unique self-adjoint realization that we denote also by $\tilde{H}_{\nu}$. The spectrum of $\tilde{H}_{\nu}$ in $L^{2, \nu}(\mathbb{D})$ consists of two parts: (i) a continuous part [ $1,+\infty[$, which corresponds to scattering states, (ii) a finite number of eigenvalues (hyperbolic Landau levels) of the form

$$
\begin{equation*}
\varepsilon_{m}^{\nu}:=4(\nu-m)(1-\nu+m), m=0,1,2, \cdots,\left[\nu-\frac{1}{2}\right] \tag{17}
\end{equation*}
$$

with infinite degeneracy, provided that $2 \nu>1$. The eigenvalues in (17) correspond eigenfunctions, which are called bound states since the particle in such a state cannot leave the system without additional energy. A concrete description of these bound states spaces will be the goal of the next section.

## 4. The bound states spaces $\mathcal{A}_{\nu, m}^{2}(\mathbb{D})$

Here, we consider the eigenspace

$$
\begin{equation*}
\mathcal{A}_{\nu, m}^{2}(\mathbb{D}):=\left\{\Phi: \mathbb{D} \rightarrow \mathbb{C}, \Phi \in L^{2, \nu}(\mathbb{D}) \text { and } \tilde{H}_{\nu} \Phi=\varepsilon_{m}^{\nu} \Phi\right\} . \tag{18}
\end{equation*}
$$

See Refs. [6, 7], for the following proposition.
Proposition 4.1. Let $2 \nu>1$ and $m=0,1,2, \cdots,\left[\nu-\frac{1}{2}\right]$.Then, we have.
(i) an orthogonal basis of $\mathcal{A}_{\nu, m}^{2}(\mathbb{D})$ is given by the functions

$$
\begin{gather*}
\phi_{k}^{\nu, m}(z):=|z|^{|m-k|}\left(1-|z|^{2}\right)^{-m} e^{-i(m-k) \arg z} \\
\times{ }_{2} F_{1}\left(-m+\frac{m-k+|m-k|}{2}, 2 \nu-m+\frac{|m-k|-m+k}{2}, 1+|m-k| ;|z|^{2}\right) \tag{19}
\end{gather*}
$$

$k=0,1,2, \cdots$, in terms of a terminating ${ }_{2} F_{1}$ Gauss hypergeometric function.
(ii) the norm square of $\phi_{k}^{\nu, m}$ in $L^{2, \nu}(\mathbb{D})$ is given by

$$
\begin{equation*}
\left\|\phi_{k}^{\nu, m}\right\|^{2}=\frac{\pi(\Gamma(1+|m-k|))^{2}}{(2(\nu-m)-1)} \frac{\Gamma\left(m-\frac{|m-k|+m-k}{2}+1\right) \Gamma\left(2 \nu-m-\frac{|m-k|+m-k}{2}\right)}{\Gamma\left(m+\frac{|m-k|-m+k}{2}+1\right) \Gamma\left(2 \nu-m+\frac{|m-k|-m+k}{2}\right)} . \tag{20}
\end{equation*}
$$

Corollary 4.1. The functions $\left\{\Phi_{k}^{\nu, m}\right\}, k=0,1,2, \ldots$, given by

$$
\begin{align*}
\Phi_{k}^{\nu, m}(z) & :=(-1)^{k}\left(\frac{2(\nu-m)-1}{\pi}\right)^{\frac{1}{2}}\left(\frac{k!\Gamma(2(\nu-m)+m)}{m!\Gamma(2(\nu-m)+k)}\right)^{\frac{1}{2}}  \tag{21}\\
& \times\left(1-|z|^{2}\right)^{-m} \bar{z}^{m-k} P_{k}^{(m-k, 2(\nu-m)-1)}(1-2 z \bar{z}), \tag{22}
\end{align*}
$$

in terms of Jacobi polynomials constitute an orthonormal basis of $\mathcal{A}_{m}^{2, \nu}(\mathbb{D})$.

Proof. Write the connection between the ${ }_{2} F_{1}$-sum and the Jacobi polynomial

$$
P_{k}^{\alpha, \beta}(u)=\frac{(1+\alpha)_{k}}{k!} \cdot{ }_{2} F_{1}\left(-k, 1+\alpha+\beta+k, 1+\alpha ; \frac{1-u}{2}\right),
$$

then the functions

$$
\begin{equation*}
\phi_{k}^{\nu, m}(z)=\frac{(-1)^{\min (m, k)}}{\left(1-|z|^{2}\right)^{m}}|z|^{m-k \mid} e^{-i(m-k) \arg z} P_{\min (m, k)}^{(|m-k|, 2(\nu-m)-1)}(1-2 z \bar{z}), \tag{23}
\end{equation*}
$$

constitute an orthonormal basis of $\mathcal{A}_{\nu, m}^{2}$. The norm square of $\phi_{k}^{\nu, m}$ in $L^{2, \nu}(\mathbb{D})$ is given by

$$
\begin{equation*}
\left\|\phi_{k}^{\nu, m}\right\|^{2}=\frac{\pi}{(2(\nu-m)-1)} \frac{(m \vee k)!\Gamma(2(\nu-m)+m \wedge k)}{(m \wedge k)!\Gamma(2(\nu-m)+m \vee k)} . \tag{24}
\end{equation*}
$$

Here, $m \wedge k:=\min (m, k)$ and $m \vee k:=\max (m, k)$. Thus, the set of functions

$$
\begin{equation*}
\Phi_{k}^{\nu, m}:=\frac{\phi_{k}^{\nu, m}}{\left\|\phi_{k}^{\nu, m}\right\|}, k=0,1,2, \ldots \tag{25}
\end{equation*}
$$

is an orthonormal basis of $\mathcal{A}_{\nu, m}^{2}(\mathbb{D})$ and can be rewritten as.

$$
\begin{align*}
\Phi_{k}^{\nu, m}(z) & =(-1)^{k}\left(\frac{2(\nu-m)-1}{\pi}\right)^{\frac{1}{2}}\left(\frac{k!\Gamma(2(\nu-m)+m)}{m!\Gamma(2(\nu-m)+k)}\right)^{\frac{1}{2}}  \tag{26}\\
& \times\left(1-|z|^{2}\right)^{-m} \bar{z}^{m-k} P_{k}^{(m-k, 2(\nu-m)-1)}(1-2 z \bar{z}) \tag{27}
\end{align*}
$$

by making appeal to the identity $(S, p .63)$ :

$$
\begin{equation*}
\frac{\Gamma(m+1)}{\Gamma(m-s+1)} P_{m}^{(-s, \alpha)}(u)=\frac{\Gamma(m+\alpha+1)}{\Gamma(m-s+\alpha+1)}\left(\frac{u-1}{2}\right)^{s} P_{m-s}^{(s, \alpha)}(u), 1 \leq s \leq m \tag{28}
\end{equation*}
$$

for $s=m-k, t=1-2|z|^{2}$, and $\alpha=2(\nu-m)-1 \ldots$
Corollary 4.2. The $L^{2}$-eigenspace $\mathcal{A}_{\nu, 0}^{2}(\mathbb{D})$, corresponding to $m=0$ in (3.1) and associated with the bottom energy $\varepsilon_{0}^{\nu}=0$ in (2.6), reduces further to the weighted Bergman space consisting of holomorphic functions $\phi: \mathbb{D} \rightarrow \mathbb{C}$ such that

$$
\begin{equation*}
\int_{\mathbb{D}}|\phi(z)|^{2}\left(1-|z|^{2}\right)^{2 \nu-2} d \mu(z)<+\infty \tag{29}
\end{equation*}
$$

## 5. Computation of the singular values $\lambda_{k}$

Elements of this basis are given in terms of Jacobi polynomials as

$$
\begin{equation*}
\phi_{k}^{\nu, m}(z)=\frac{(-1)^{\min (m, k)}}{\left(1-|z|^{2}\right)^{m}}|z|^{|m-k|} e^{-i(m-k) \operatorname{argz}} P_{\min (m, k)}^{(|m-k|, 2(\nu-m)-1)}(1-2 z \bar{z}) . \tag{30}
\end{equation*}
$$

The norm square of $\phi_{k}^{\nu, m}$ in $L^{2, \nu}(\mathbb{D})$ is given by

$$
\begin{equation*}
\rho_{k}^{\nu, m}=\frac{\pi}{(2(\nu-m)-1)} \frac{(m \vee k)!\Gamma(2(\nu-m)+m \wedge k)}{(m \wedge k)!\Gamma(2(\nu-m)+m \vee k)} \tag{31}
\end{equation*}
$$

Here, $m \wedge k:=\min (m, k)$ and $m \vee k:=\max (m, k)$. Let us introduce the notation. The set of functions

$$
\begin{equation*}
\gamma_{k}^{\nu, m}:=\frac{(-1)^{m \wedge k}}{\sqrt{\rho_{k}^{\nu, m}}}, k=0,1,2, \ldots \tag{32}
\end{equation*}
$$

So that we consider the elements

$$
\begin{equation*}
\Phi_{k}^{\nu, m}(z):=\gamma_{k}^{\nu, m} \frac{1}{(1-z \bar{z})^{m}}|z|^{|m-k|} e^{-i(m-k) \arg z} P_{\min (m, k)}^{(|m-k|, 2(\nu-m)-1)}(1-2 z \bar{z}) . \tag{33}
\end{equation*}
$$

### 5.1 The action $\mathcal{L}_{\nu}$

Lemma 5.1. We set $z=\rho e^{i t}$, and $I=-\int_{0}^{2 \pi} e^{i(k-m) \theta} \log \left(\left|z-r e^{i \theta}\right|\right) \frac{d \theta}{2 \pi}$, we have

$$
\begin{cases}I=-\log (\rho \wedge r) & k=m,  \tag{34}\\ I=\frac{e^{i(k-m) t}}{2|m-k|}\left(\left(\frac{r}{\rho}\right)^{m-k} \wedge\left(\frac{r}{\rho}\right)^{m-k}\right) & k \neq m,\end{cases}
$$

Proof. By ref. [3], it remains to prove that this lemma for $k<m$. We have

$$
\begin{equation*}
\int_{0}^{2 \pi} e^{i(k-m) \theta} \log \left(\left|\rho e^{i t}-r e^{i \theta}\right|\right) d \theta=-\int_{0}^{2 \pi} e^{i(m-k)(-\theta)} \log \left(\left|r e^{i(-t)}-\rho e^{i(-\theta)}\right|\right) d(-\theta) \tag{35}
\end{equation*}
$$

The function $\theta \rightarrow e^{i(m-k)(-\theta)} \log \left(\left|r e^{i(-t)}-\rho e^{i(-\theta)}\right|\right)$ is a periodic mapping with the period equal $2 \pi$, then

$$
\begin{gathered}
\int_{0}^{2 \pi} e^{i(k-m) \theta} \log \left(\left|\rho e^{i t}-r e^{i \theta}\right|\right) d \theta \\
=-\int_{0}^{2 \pi} e^{i(m-k)(-\theta)} \log \left(\left|r e^{i(-t)}-\rho e^{i(-\theta)}\right|\right) d(-\theta) \\
=\frac{e^{i(k-m) t}}{2(m-k)} \times\left(\left(\frac{r}{\rho}\right)^{m-k} \wedge\left(\frac{r}{\rho}\right)^{m-k}\right)
\end{gathered}
$$

Lemma 5.2. For all $\lambda \in \partial D$. $\mathcal{L}_{\nu}$ commutes with the rotations $R_{\lambda}$, where

$$
\left(R_{\lambda} f\right)(z)=f(\lambda z) .
$$

Proof. We observe that

$$
R_{\lambda} \phi_{k}^{\nu, m}(z)=\lambda^{k-m} \phi_{k}^{\nu, m}(z), \quad \forall k \neq m .
$$

Corollary 5.1. $\left\{\mathcal{L}_{\nu}\left(\phi_{k}^{\nu, m}\right)\right\}_{k=0}^{\infty}$ are orthonormal in $L^{2, \nu}(\mathbb{D})$.
Proof. As $R_{\lambda}$ is an isometry of $L^{2, \nu}(\mathbb{D})$,

$$
\begin{gathered}
\left(\mathcal{L}_{\nu}\left(\phi_{k}^{\nu, m}\right), \mathcal{L}_{\nu}\left(\phi_{j}^{\nu, m}\right)\right) \\
=\left(R_{\lambda} \mathcal{L}_{\nu}\left(\phi_{k}^{\nu, m}\right), R_{\lambda} \mathcal{L}_{\nu}\left(\phi_{j}^{\nu, m}\right)\right) \\
=\left(\mathcal{L}_{\nu} R_{\lambda}\left(\phi_{k}^{\nu, m}\right), \mathcal{L}_{\nu} R_{\lambda}\left(\phi_{j}^{\nu, m}\right)\right) \\
=\overline{\lambda^{j-k}}\left(\mathcal{L}_{\nu}\left(\phi_{k}^{\nu, m}\right), \mathcal{L}_{\nu}\left(\phi_{j}^{\nu, m}\right)\right), \text { if } m>k,
\end{gathered}
$$

or

$$
=\lambda^{k-j}\left(\mathcal{L}_{\nu}\left(\phi_{k}^{\nu, m}\right), \mathcal{L}_{\nu}\left(\phi_{j}^{\nu, m}\right)\right), \text { if } m<k
$$

For all $\lambda \in \partial D$, since $\lambda \neq 0$, we have

$$
\left(\mathcal{L}_{\nu}\left(\phi_{k}^{\nu, m}\right), \mathcal{L}_{\nu}\left(\phi_{j}^{\nu, m}\right)\right)=0 \text { if } j \neq k
$$

Lemma 5.3. If we denote ${ }^{1} \phi_{k}^{\nu, m}(z)$, if $k>m$ and ${ }^{2} \phi_{k}^{\nu, m}(z)$, if $k<m$, we have

$$
\mathcal{L}_{\nu}\left({ }^{1} \phi_{k}^{\nu, m}\right)(z)=\frac{\Gamma(k+1) \Gamma(2 \nu-m)}{\Gamma(m+1) \Gamma(2 \nu-k)} \mathcal{L}_{\nu}\left({ }^{2} \phi_{k}^{\nu, m}\right)(z) .
$$

Proof. Just use

$$
\begin{equation*}
\frac{\Gamma(m+1)}{\Gamma(m-s+1)} P_{m}^{(-s, \alpha)}(u)=\frac{\Gamma(m+\alpha+1)}{\Gamma(m-s+\alpha+1)}\left(\frac{u-1}{2}\right)^{s} P_{m-s}^{(s, \alpha)}(u), 1 \leq s \leq m . \tag{36}
\end{equation*}
$$

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Proposition 5.1. The action of the operator $\mathcal{L}$ on a basis element $\phi_{k}^{\nu, m}$ is of the form: If $k=m$, We put $z=\rho e^{i \theta}$ then
$\mathcal{L}_{\nu}\left(\phi_{k}^{\nu, m}\right)(z)=\frac{\alpha_{k}^{\nu, m}}{2(2 \nu-m+1)} \sqrt{\frac{2(\nu-m)-1}{\pi}}\left(1-\rho^{2}\right)^{2 \nu-m-1}{ }_{3} F_{2}\left(\left.\begin{array}{c}-m+1,2 \nu-m, 2 \nu-m+1 \\ 2(\nu-m), 2 \nu-m+2\end{array} \right\rvert\, 1-\rho^{2}\right)$.

If $k \neq m$ then

$$
\mathcal{L}_{\nu}\left(\phi_{k}^{\nu, m}\right)(z)=\frac{\pi \gamma_{k}^{\nu, m} e^{i(k-m) t}}{2(k-m)}\left(I_{3}+I_{4}\right)
$$

where

$$
I_{3}=\frac{(1+k-m)_{m}}{m!(k-m+1)} \rho^{k-m+2}\left(1-\rho^{2}\right)^{2 \nu-m-1}{ }_{2} F_{1}\left(\begin{array}{cc}
-m+1,2(\nu-m)+k & \mid \rho^{2} \\
2+k-m &
\end{array}\right),
$$

and

$$
I_{4}=\frac{\alpha_{k}^{\nu, m}}{2 \nu-m-1}\left(1-\rho^{2}\right)^{2 \nu-m-1}{ }_{2} F_{1}\left(\begin{array}{cc}
-m+1,2 \nu-m-1 & \mid \rho^{2} \\
2(\nu-m),
\end{array}\right) .
$$

Proof. For $k=m$, we have

$$
\begin{aligned}
\mathcal{L}_{\nu}\left(\phi_{k}^{\nu, m}\right)(z) & =\frac{(-1)^{m}}{\pi} \sqrt{\frac{2(\nu-m)-1}{\pi}} \int_{\mathbb{D}}\left(1-|\xi|^{2}\right)^{2 \nu-m-2} P_{m}^{(0,2(\nu-m)-1)}\left(1-2|\xi|^{2}\right) \log (|z-\xi|) d \mu(\xi) \\
& =(-1)^{m} \sqrt{\frac{2(\nu-m)-1}{\pi}} \int_{0}^{1}\left(1-r^{2}\right)^{2 \nu-m-2} P_{m}^{(0,2(\nu-m)-1)}\left(1-2 r^{2}\right) \log (\rho \wedge r) d r^{2} \\
& =\frac{(-1)^{m}}{2} \sqrt{\frac{2(\nu-m)-1}{\pi}} \int_{0}^{1}(1-t)^{2 \nu-m-2} P_{m}^{(0,2(\nu-m)-1)}(1-2 t) \log \left(\rho^{2} \vee t\right) d t \\
& =\frac{(-1)^{m}}{2} \sqrt{\frac{2(\nu-m)-1}{\pi}}\left[I_{1}+I_{2}\right] .
\end{aligned}
$$

where

$$
I_{1}=\int_{0}^{\rho^{2}}(1-t)^{2 \nu-m-2} P_{m}^{(0,2(\nu-m)-1)}(1-2 t) \log \left(\rho^{2} \vee t\right) d t
$$

and

$$
I_{2}=\int_{\rho^{2}}^{1}(1-t)^{2 \nu-m-2} P_{m}^{(0,2(\nu-m)-1)}(1-2 t) \log (t) d t .
$$

Calculus of $I_{1}$.

$$
I_{1}=\log \left(\rho^{2}\right) \int_{\rho^{2}}^{1}(1-t)^{2 \nu-m-2} P_{m}^{(0,2(\nu-m)-1)}(1-2 t) d t .
$$

We use the formula

$$
P_{k}^{(\alpha, \beta)}(u)=\frac{(1+\alpha)_{k}}{k!} F_{1}\left(\begin{array}{cc}
-k, 1+\alpha+\beta+k & 1-u \\
1+\alpha & \left\lvert\, \frac{1-\alpha}{2}\right.
\end{array}\right) .
$$

We have

$$
I_{1}=\log \left(\rho^{2}\right) \int_{0}^{\rho^{2}}(1-t)^{2 \nu-m-2} F_{1}\left(\begin{array}{cc}
-m, 2 \nu-m & \mid t \\
1 & \mid t .
\end{array}\right.
$$

By Ref. [8], we have

$$
\int x^{c-1}(1-x)^{b-c-1}{ }_{2} F_{1}\left(\begin{array}{cc}
a, b & \mid x \\
c &
\end{array}\right) d x=\frac{1}{c} x^{c}(1-x)^{b-c} F_{1}\left(\begin{array}{cc}
a+1, b & \mid x \\
c+1 &
\end{array}\right)
$$

implies that

$$
I_{1}=\log \left(\rho^{2}\right) \rho^{2}\left(1-\rho^{2}\right)^{2 \nu-m-1} F_{1}\left(\begin{array}{cc}
-m+1,2 \nu-m & \mid \rho^{2} \\
2 &
\end{array}\right)
$$

Calculus of $I_{2}$.

$$
I_{2}=\int_{\rho^{2}}^{1}(1-t)^{2 \nu-m-2} P_{m}^{(0,2(\nu-m)-1)}(1-2 t) \log (t) d t
$$

Use the previous formula in Ref. [8] and the integration by part gives

$$
\left.\begin{array}{rl}
I_{2} & =\left[t 1-t^{2 \nu-m-1} F_{2}\binom{-m+1,2 \nu-m}{2} \log t\right]_{\rho^{2}}^{1}-\int_{\rho^{2}}^{1} 1-t^{2 \nu-m}{ }_{2} F_{1}\binom{-m+1,2 \nu-m}{2} d t \\
& =-\rho^{2} \log \rho^{2} 1-\rho^{2 \nu-m-1}{ }_{2} F_{1}\left(\left.\begin{array}{c}
-m+1,2 \nu-m \\
2
\end{array} \right\rvert\, \rho^{2}\right.
\end{array}\right)-\int_{\rho^{2}}^{1} 1-t^{2 \nu-m}{ }_{2} F_{1}\binom{-m+1,2 \nu-m}{2} d t .
$$

Calculus of

$$
\int_{\rho^{2}}^{1}(1-t)^{2 \nu-m}{ }_{2} F_{1}\left(\begin{array}{cc}
-m+1,2 \nu-m \\
2 & \mid t
\end{array}\right) d t .
$$

Use the following formula, which has place in [9]

$$
\begin{aligned}
{ }_{2} F_{1}\left(\left.\begin{array}{c}
a, b \\
c
\end{array} \right\rvert\, t\right)= & \frac{\Gamma(c) \Gamma(c-a-b)}{\Gamma(c-a) \Gamma(c-b)} F_{2}\left(\left.\begin{array}{c}
a, b \\
a+b-c+1
\end{array} \right\rvert\, 1-t\right) \\
& +\frac{\Gamma(c) \Gamma(a+b-c)}{\Gamma(a) \Gamma(b)}(1-t)^{c-a-b}{ }_{2} F_{1}\left(\left.\begin{array}{c}
a, b \\
a+b-c+1
\end{array} \right\rvert\, 1-t\right)
\end{aligned}
$$

We put $a=1-m, b=2 \nu-m, c=2$ and use the formula Boher-Mollerup, for $z \in \mathrm{IR}_{+}^{*}$,

$$
\Gamma(z)=\frac{e^{-\gamma z}}{z} \prod_{n=1}^{\infty}\left(1+\frac{z}{n}\right)^{-1} e^{-\frac{z}{n}},
$$

which implies $\frac{1}{\Gamma(1-m)}=0$, then

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$$
{ }_{2} F_{1}\left(\begin{array}{cc}
-m+1,2 \nu-m & \mid t \\
2 &
\end{array}\right)=\frac{2 \Gamma(2(m-\nu)+1)}{m!\Gamma(2+m-2 \nu)_{2}} F_{1}\left(\left.\begin{array}{c}
-m+1,2 \nu-m \\
2(\nu-m)
\end{array} \right\rvert\, 1-t\right)
$$

implies that

$$
\int_{\rho^{2}}^{1}(1-t)^{2 \nu-m}{ }_{2} F_{1}\left(\left.\begin{array}{c}
-m+1,2 \nu-m \\
2
\end{array} \right\rvert\, t\right) d t=\frac{2 \Gamma(2(m-\nu)+1)}{m!\Gamma(2+m-2 \nu)} \int_{\rho^{2}}^{1}(1-t)^{2 \nu-m}{ }_{2} F_{1}\left(\left.\begin{array}{c}
-m+1,2 \nu-m \\
2(\nu-m)
\end{array} \right\rvert\, 1-t\right) d t .
$$

By the change $1-t=s$, we get

$$
\int_{\rho^{2}}^{1}(1-t)^{2 \nu-m} F_{1}\left(\left.\begin{array}{c}
-m+1,2 \nu-m \\
2
\end{array} \right\rvert\, t\right) d t=\frac{2 \Gamma(2(m-\nu)+1)}{m!\Gamma(2+m-2 \nu)} \int_{0}^{1-\rho^{2}} t^{2 \nu-m} F_{1}\left(\left.\begin{array}{c}
-m+1,2 \nu-m \\
2(\nu-m)
\end{array} \right\rvert\, t\right) d t .
$$

In [8], p. 44,

$$
\int x^{\alpha-1}{ }_{2} F_{1}\left(\begin{array}{cc}
a, b & \mid-t \\
c & \left\lvert\,-\alpha=\frac{x^{\alpha}}{\alpha_{3}} F_{2}\left(\left.\begin{array}{c}
a, b, \alpha \\
c, \alpha+1
\end{array} \right\rvert\,-t\right)+\frac{\Gamma(\alpha) \Gamma(a-\alpha) \Gamma(b-\alpha) \Gamma(c)}{\Gamma(a) \Gamma(b) \Gamma(c-\alpha)}\right., ~
\end{array}\right.
$$

Since $a=1-m, b=2 \nu-m, c=2(\nu-m)$, and $\alpha=2 \nu-m+1$ we have

$$
\frac{\Gamma(\alpha) \Gamma(a-\alpha) \Gamma(b-\alpha) \Gamma(c)}{\Gamma(a) \Gamma(b) \Gamma(c-\alpha)}=0
$$

and by the change $t=-s$

$$
\begin{gathered}
\int_{0}^{1-\rho^{2}} t^{2 \nu-m+1}{ }_{2} F_{1}\left(\begin{array}{c}
-m+1,2 \nu-m \\
2(\nu-m)
\end{array}|t|\right) d t \\
=(-1)^{m} \int_{0}^{\rho^{2}} t^{2 \nu-m}{ }_{2} F_{1}\left(\begin{array}{c}
-m+1,2 \nu-m \\
2(\nu-m)
\end{array}|-t|\right) d t \\
=(-1)^{m} \frac{\left(\rho^{2}-1\right)^{2 \nu-m+1}}{2 \nu-m+1} F_{3} F_{2}\left(\begin{array}{c}
-m+1,2 \nu-m, 2 \nu-m+1 \\
2(\nu-m), 2 \nu-m+2
\end{array}\left|1-\rho^{2}\right|\right)
\end{gathered}
$$

we set $\alpha_{k}^{\nu, m}=\frac{2 \Gamma(2(m-\nu)+1)}{m!\Gamma(2+m-2 \nu)}$. We get

$$
\begin{aligned}
I_{2}= & -\rho^{2} \log \rho^{2} 1-\rho^{2 \nu-m-1}{ }_{2} F_{1}\left(\left.\begin{array}{c}
-m+1,2 \nu-m \\
2
\end{array} \right\rvert\, \rho^{2}\right) \\
& +-1^{m} \alpha_{k}^{\nu, m \frac{1-\rho^{2 \nu-m-1}}{2 \nu-m+1}{ }_{3}} F_{2}\left(\left.\begin{array}{c}
-m+1,2 \nu-m, 2 \nu-m+1 \\
2 \nu-m, 2 \nu-m+2
\end{array} \right\rvert\, 1-\rho^{2}\right) .
\end{aligned}
$$

Finally

$$
\mathcal{L}_{\nu}\left(\phi_{k}^{\nu, m}\right)(z)=\frac{\alpha_{k}^{\nu, m}}{2(2 \nu-m+1)} \sqrt{\frac{2(\nu-m)-1}{\pi}}\left(1-\rho^{2}\right)^{2 \nu-m-1}{ }_{3} F_{2}\left(\left.\begin{array}{c}
-m+1,2 \nu-m, 2 \nu-m+1 \\
2(\nu-m), 2 \nu-m+2
\end{array} \right\rvert\, 1-\rho^{2}\right) .
$$

Now if $k>m$, set $z=\rho e^{i t}$.

$$
\begin{aligned}
& \mathcal{L}_{\nu}\left(\phi_{k}^{\nu, m}\right)(z)=\gamma_{k}^{\nu, m} \int_{\mathbb{D}}\left(1-|\xi|^{2}\right)^{2 \nu-m-2} \xi^{k-m} \log \left(\frac{1}{|z-\xi|}\right) P_{m}^{(k-m, 2(\nu-m)-1)}\left(1-2|\xi|^{2}\right) d \mu(\xi) \\
& =\gamma_{k}^{\nu, m} \int_{0}^{1}\left(1-r^{2}\right)^{2 \nu-m-2} r^{k-m+1} P_{m}^{(k-m, 2(\nu-m)-1)}\left(1-2 r^{2}\right) \int_{0}^{2 \pi} e^{i(k-m) \theta} \log \left(\frac{1}{\left|z-r^{i \theta}\right|}\right) d \theta d r \\
& =\frac{\pi \gamma_{k}^{\nu, m} e^{i(k-m) t}}{2(k-m)} \int_{0}^{1}\left(1-r^{2}\right)^{2 \nu-m-2} r^{k-m} P_{m}^{(k-m, 2(\nu-m)-1)}\left(1-2 r^{2}\right)\left(\left(\frac{r}{\rho}\right)^{k-m} \wedge\left(\frac{\rho}{r}\right)^{k-m}\right) d r^{2} \\
& =\frac{\pi \gamma_{k}^{\nu, m} e^{i(k-m) t}}{2(k-m)}\left(\int_{0}^{\rho}\left(1-r^{2}\right)^{2 \nu-m-2} r^{k-m} P_{m}^{(k-m, 2(\nu-m)-1)}\left(1-2 r^{2}\right)\left(\left(\frac{r}{\rho}\right)^{k-m} \wedge\left(\frac{\rho}{r}\right)^{k-m}\right) d r^{2}\right. \\
& \quad+\int_{\rho}^{1}\left(1-r^{2}\right)^{2 \nu-m-2} r^{k-m} P_{m}^{(k-m, 2(\nu-m)-1)}\left(1-2 r^{2}\right)\left(\left(\frac{r}{\rho}\right)^{k-m} \wedge\left(\frac{\rho}{r}\right)^{k-m}\right) d r^{2} .
\end{aligned}
$$

We set

$$
I_{3}=\int_{0}^{\rho}\left(1-r^{2}\right)^{2 \nu-m-2} r^{k-m} P_{m}^{(k-m, 2(\nu-m)-1)}\left(1-2 r^{2}\right)\left(\left(\frac{r}{\rho}\right)^{k-m} \wedge\left(\frac{\rho}{r}\right)^{k-m}\right) d r^{2}
$$

and

$$
I_{4}=\int_{\rho}^{1}\left(1-r^{2}\right)^{2 \nu-m-2} r^{k-m} P_{m}^{(k-m, 2(\nu-m)-1)}\left(1-2 r^{2}\right)\left(\left(\frac{r}{\rho}\right)^{k-m} \wedge\left(\frac{\rho}{r}\right)^{k-m}\right) d r^{2}
$$

Calculus of $I_{3}$.

$$
I_{3}=\frac{\rho^{m-k}(1+k-m)_{m}}{m!} \int_{0}^{\rho^{2}} t^{k-m}(1-t)^{2 \nu-m-2}{ }_{2} F_{1}\left(\begin{array}{cc}
-m, 2(\nu-m)+k & \mid t \\
1+k-m &
\end{array}\right) d t .
$$

By the formula

$$
\int x^{c-1}(1-x)^{b-c-1}{ }_{2} F_{1}\left(\begin{array}{cc}
a, b & \mid x \\
c &
\end{array}\right) d x=\frac{1}{c} x^{c}(1-x)^{b-c} F_{1}\left(\begin{array}{cc}
a+1, b & \mid x \\
c+1 &
\end{array}\right),
$$

we have

$$
I_{3}=\frac{(1+k-m)_{m}}{m!(k-m+1)} \rho^{k-m+2}\left(1-\rho^{2}\right)^{2 \nu-m-1}{ }_{2} F_{1}\left(\begin{array}{cc}
-m+1,2(\nu-m)+k & \mid \rho^{2} \\
2+k-m &
\end{array}\right)
$$

Calculus of $I_{4}$.

$$
\begin{aligned}
I_{4} & =\int_{\rho}^{1}\left(1-r^{2}\right)^{2 \nu-m-2} r^{k-m} P_{m}^{(k-m, 2(\nu-m)-1)}\left(1-2 r^{2}\right)\left(\left(\frac{r}{\rho}\right)^{k-m} \wedge\left(\frac{\rho}{r}\right)^{k-m}\right) d r^{2} \\
& =\frac{\rho^{k-m}(1+k-m)_{m}}{2 m!} \int_{\rho^{2}}^{1}(1-t)^{2 \nu-m-2}{ }_{2} F_{1}\left(\begin{array}{cc}
-m, 2(\nu-m)+k & \mid t) d t . \\
1+k-m &
\end{array} .\right.
\end{aligned}
$$

As the previous

$$
\begin{gathered}
\int_{\rho^{2}}^{1}(1-t)^{2 \nu-m-2}{ }_{2} F_{1}\left(\left.\begin{array}{c}
-m, 2(\nu-m)+k \\
1+k-m
\end{array} \right\rvert\, t\right) d t=\alpha_{k}^{\nu, m} \int_{\rho^{2}}^{1}(1-t)^{2 \nu-m-2}{ }_{2} F_{1}\left(\left.\begin{array}{c}
-m+1,2 \nu-m \\
2(\nu-m)
\end{array} \right\rvert\, 1-t\right) d t \\
=(-1)^{m} \alpha_{k}^{\nu, m} \int_{\rho^{2}-1}^{0} t^{2 \nu-m-2}{ }_{2} F_{1}\left(\left.\begin{array}{c}
-m+1,2 \nu-m \\
2(\nu-m)
\end{array} \right\rvert\, t\right) d t \\
=\frac{\alpha_{k}^{\nu, m}}{2 \nu-m-1}\left(1-\rho^{2}\right)^{2 \nu-m-1}{ }_{3} F_{2}\left(\left.\begin{array}{c}
-m+1,2 \nu-m, 2 \nu-m-1 \\
2(\nu-m), 2 \nu-m
\end{array} \right\rvert\, 1-\rho^{2}\right)
\end{gathered}
$$

also

$$
{ }_{3} F_{2}\left(\begin{array}{cc}
-m+1,2 \nu-m, 2 \nu-m-1 & \mid 1-\rho^{2} \\
2(\nu-m), 2 \nu-m &
\end{array}\right)={ }_{2} F_{1}\left(\begin{array}{cc}
-m+1,2 \nu-m-1 & \mid 1-\rho^{2} \\
2(\nu-m) &
\end{array}\right)
$$

Now if $k<m$. We have
$\phi_{k}^{\nu, m}(z)=(-1)^{k} \sqrt{\frac{2(\nu-m)-1}{\pi} \frac{k!\Gamma(2(\nu-m)+m)}{m!\Gamma(2(\nu-m)+k)}}\left(1-|z|^{2}\right)^{-m} z^{m-k} P_{k}^{(m-k, 2(\nu-m)-1)}\left(1-2|z|^{2}\right)$.
By the formula

$$
\begin{equation*}
\frac{\Gamma(m+1)}{\Gamma(m-s+1)} P_{m}^{(-s, \alpha)}(u)=\frac{\Gamma(m+\alpha+1)}{\Gamma(m-s+\alpha+1)}\left(\frac{u-1}{2}\right)^{s} P_{m-s}^{(s, \alpha)}(u), 1 \leq s \leq m \tag{38}
\end{equation*}
$$

and put $s=m-k$ and $\alpha=2(\nu-m)-1$, we have

$$
P_{k}^{(m-k, 2(\nu-m)-1)}\left(1-2|z|^{2}\right)=\frac{m!\Gamma(k+\alpha+1)}{k!\Gamma(m+\alpha+1)} P_{m}^{(k-m, 2(\nu-m)-1)}\left(1-2|z|^{2}\right)
$$

substituting in the expression of $\phi_{k}^{\nu, m}(z)$, we get

$$
\phi_{k}^{\nu, m}(z)=(-1)^{m} \sqrt{\frac{2(\nu-m)-1}{\pi} \frac{m!\Gamma(2(\nu-m)+k)}{k!\Gamma(2(\nu-m)+m)}}\left(1-|z|^{2}\right)^{-m} z^{k-m} P_{m}^{(k-m, 2(\nu-m)-1)}\left(1-2|z|^{2}\right),
$$

it is the same formula for $k>m$, which proves the same formula of $\mathcal{L}_{\nu}\left(\phi_{k}^{\nu, m}\right)(z)$ if $k>m$.

Remark 5.1. By the previous formula in [9], we have

$$
{ }_{2} F_{1}\left(\begin{array}{cc}
-m+1,2(\nu-m)+k & \mid \rho^{2} \\
2(\nu-m) & k!\Gamma(2+k-m) \\
\Gamma(1-2(\nu-m))_{2} & F_{1}\left(\left.\begin{array}{c}
-m+1,2(\nu-m)+k \\
2(\nu-m)
\end{array} \right\rvert\, 1-\rho^{2}\right.
\end{array}\right) .
$$

### 5.2 The spectrum of $\mathcal{L}_{\nu}$

Proposition 5.2. If $k \neq m$, then

$$
\lambda_{k}=\sqrt{J_{1}+J_{2}+J_{3}} .
$$

where

$$
\begin{gathered}
J_{1}=\left(\frac{(1+k-m)_{m}}{m!(k-m+1)}\right)^{2} \sum_{n=0}^{\infty} A_{n} \frac{\Gamma(2 n+2 k-2 m+6-1) \Gamma(4 \nu-2 m-1)}{\Gamma(2 n+2 k-4 m+4 \nu+6)} \\
J_{2}=\left(\frac{\alpha_{k}^{\nu, m}}{2 \nu-m-1}\right)^{2} \sum_{n=0}^{\infty} A_{n} \frac{\Gamma(4 \nu-2 m-1) \Gamma(2 n+2)}{\Gamma(2 n+4 \nu-2 m+1)}
\end{gathered}
$$

and

$$
J_{3}=\frac{(1+k-m)_{m} \alpha_{k}^{\nu, m}}{m!(k-m+1)(2 \nu-m-1)}\left(\sum_{n=0}^{\infty} A_{n} \frac{\Gamma(k-m+2) \Gamma(4 \nu-2 m-1)}{\Gamma(4 \nu-k-3 m)}\right)
$$

If $k=m$ then

$$
\begin{equation*}
\lambda_{k}^{2}=\frac{\alpha_{k}^{\nu, m}(2(2 \nu-m)-1)}{8(\pi(2 \nu-m+1))} \sum_{n=0}^{\infty} \frac{B_{n}}{n+2 \nu-m} . \tag{39}
\end{equation*}
$$

where

$$
B_{n}=\sum_{n=0}^{\infty} \frac{\Gamma(-m+1) \Gamma(2 \nu-m) \Gamma(2(\nu-m)+1)}{n!\Gamma(2(\nu-m) \Gamma(2 \nu-m+2)} .
$$

Proof. If $k \neq m$. We have

$$
\left(\mathcal{L}_{\nu}\left(\phi_{k}^{\nu, m}\right)\right)(z)=\frac{\pi \gamma_{k}^{\nu, m}\left(I_{3}+I_{4}\right)}{2(k-m)} e^{i(k-m) t}
$$

We set $H=\left(L^{2}(\mathbb{D}),\left(1-|\xi|^{2}\right)^{2 \nu-2} d \mu(\xi)\right), I_{3}=I_{3}(\rho)$, and $I_{4}=I_{4}(\rho)$ we have

$$
\begin{gathered}
\lambda_{k}^{2}=\left\langle\mathcal{L}_{\nu}\left(\phi_{k}^{\nu, m}\right), \mathcal{L}_{\nu}\left(\phi_{k}^{\nu, m}\right)\right\rangle_{H} \\
=\frac{\pi^{2} \gamma_{k}^{\nu, m}}{(k-m)} \int_{0}^{1}\left(I_{3}(\rho)+I_{4}(\rho)\right)^{2} \rho d \rho .
\end{gathered}
$$

Calculus of $\int_{0}^{1}\left(I_{3}(\rho)\right)^{2} \rho d \rho$.

$$
I_{3}(\rho)=\frac{(1+k-m)_{m}}{m!(k-m+1)} \rho^{k-m+2}\left(1-\rho^{2}\right)^{2 \nu-m-1}{ }_{2} F_{1}\left(\begin{array}{cc}
-m+1,2(\nu-m)+k & \mid \rho^{2} \\
2+k-m &
\end{array}\right.
$$

Since

$$
{ }_{2} F_{1}\left(\left.\begin{array}{cc}
-m+1,2(\nu-m)+k \\
2+k-m
\end{array} \right\rvert\, \rho^{2}\right)=\sum_{n=0}^{\infty} \frac{(-m+1)_{n}(2(\nu-m)+k)_{n}}{(2+k-m)_{n}} \frac{\rho^{2 n}}{n!},
$$

then

$$
\left(I_{3}(\rho)\right)^{2}=\left(\frac{(1+k-m)_{m}}{m!(k-m+1)}\right)^{2} \sum_{n=0}^{\infty} A_{n} \rho^{2 n}\left(1-\rho^{2}\right)^{4 \nu-2 m-2}
$$

where

$$
A_{n}=\frac{1}{n!} \sum_{i=0}^{n} \frac{(-m+1)_{i}(-m+1)_{n-i}(2(\nu-m)+k)_{i}(2(\nu-m)+k)_{n-i}}{(2(\nu-m))_{i}(2(\nu-m))_{n-i}} .
$$

Thus

$$
J_{1}=\int_{0}^{1}\left(I_{3}(\rho)\right)^{2} \rho d \rho=\left(\frac{(1+k-m)_{m}}{m!(k-m+1)}\right)^{2} \sum_{n=0}^{\infty} A_{n} \int_{0}^{1} \rho^{2 n+2 k-2 m+6-1}\left(1-\rho^{2}\right)^{4 \nu-2 m-1-1} d \rho .
$$

Use the fact that

$$
\int_{0}^{1} t^{\alpha-1}(1-t)^{\beta-1} d t=\frac{\Gamma(\alpha) \Gamma(\beta)}{\Gamma(\alpha+\beta)}
$$

implies

$$
\begin{equation*}
\int_{0}^{1}\left(I_{3}(\rho)\right)^{2} \rho d \rho=\left(\frac{(1+k-m)_{m}}{m!(k-m+1)}\right)^{2} \sum_{n=0}^{\infty} A_{n} \frac{\Gamma(2 n+2 k-2 m+6-1) \Gamma(4 \nu-2 m-1)}{\Gamma(2 n+2 k-4 m+4 \nu+6)} . \tag{40}
\end{equation*}
$$

Calculus of $\int_{0}^{1}\left(I_{4}(\rho)\right)^{2} \rho d \rho$.
In the same,

$$
\begin{equation*}
J_{2}=\int_{0}^{1}\left(I_{4}(\rho)\right)^{2} \rho d \rho=\left(\frac{\alpha_{k}^{\nu, m}}{2 \nu-m-1}\right)^{2} \sum_{n=0}^{\infty} A_{n} \frac{\Gamma(4 \nu-2 m-1) \Gamma(2 n+2)}{\Gamma(2 n+4 \nu-2 m+1)} . \tag{41}
\end{equation*}
$$

Calculus of $2 \int_{0}^{1}\left(I_{3}(\rho)\right)\left(I_{4}(\rho)\right) \rho d \rho$.
$J_{3}=2 \int_{0}^{1}\left(I_{3}(\rho)\right)\left(I_{4}(\rho)\right) \rho d \rho=\frac{(1+k-m)_{m} \alpha_{k}^{\nu, m}}{m!(k-m+1)(2 \nu-m-1)} \sum_{n=0}^{\infty} A_{n} \frac{\Gamma(k-m+2) \Gamma(4 \nu-2 m-1)}{\Gamma(4 \nu-k-3 m)}$.

If $k=m$.
Since

$$
\begin{align*}
& \left({ }_{3} F_{2}\left(\left.\begin{array}{c}
-m+1,2 \nu-m, 2(\nu-m)+1 \\
2(\nu-m), 2 \nu-m+2
\end{array} \right\rvert\, 1-\rho^{2}\right)\right)^{2}=\sum_{n=0}^{\infty} \frac{\Gamma(-m+1) \Gamma(2 \nu-m) \Gamma(2(\nu-m)+1)}{n!\Gamma(2(\nu-m) \Gamma(2 \nu-m+2)}\left(1-\rho^{2}\right)^{n} .  \tag{43}\\
& \lambda_{k}^{2}=\frac{\alpha_{k}^{\nu, m}(2(2 \nu-m)-1)}{8(\pi(2 \nu-m+1))} \sum_{n=0}^{\infty} B_{n} \int_{0}^{1}\left(1-\rho^{2}\right)^{n+2 \nu-m-1} d \rho=\frac{\alpha_{k}^{\nu, m}(2(2 \nu-m)-1)}{8(\pi(2 \nu-m+1))} \sum_{n=0}^{\infty} \frac{B_{n}}{n+2 \nu-m}, \tag{44}
\end{align*}
$$

where

$$
B_{n}=\sum_{n=0}^{\infty} \frac{\Gamma(-m+1) \Gamma(2 \nu-m) \Gamma(2(\nu-m)+1)}{n!\Gamma(2(\nu-m) \Gamma(2 \nu-m+2)} .
$$

## 6. Asymptotic behavior of singular values $\lambda_{\boldsymbol{k}}$ as $\boldsymbol{k} \rightarrow \infty$

## Proposition 6.1.

$$
\lambda_{k} \sim C \sqrt{k^{m-4 \nu+1}}, \text { as } k \rightarrow \infty
$$

where $C$ is a constant.
Proof. If $k>m$, then

$$
\lambda_{k}=\sqrt{J_{1}+J_{2}+J_{3}},
$$

where

$$
\begin{gathered}
J_{1}=\left(\frac{(1+k-m)_{m}}{m!(k-m+1)}\right)^{2} \sum_{n=0}^{\infty} A_{n} \frac{\Gamma(2 n+2 k-2 m+6-1) \Gamma(4 \nu-2 m-1)}{\Gamma(2 n+2 k-4 m+4 \nu+6)} \\
J_{2}=\left(\frac{\alpha_{k}^{\nu, m}}{2 \nu-m-1}\right)^{2} \sum_{n=0}^{\infty} A_{n} \frac{\Gamma(4 \nu-2 m-1) \Gamma(2 n+2)}{\Gamma(2 n+4 \nu-2 m+1)}
\end{gathered}
$$

and

$$
J_{3}=\frac{(1+k-m)_{m} \alpha_{k}^{\nu, m}}{m!(k-m+1)(2 \nu-m-1)}\left(\sum_{n=0}^{\infty} A_{n} \frac{\Gamma(k-m+2) \Gamma(4 \nu-2 m-1)}{\Gamma(4 \nu-k-3 m)}\right)
$$

The limit of $\lambda_{k}$ as $k \rightarrow \infty$.
We use the formula

$$
\frac{\Gamma(k+a)}{k=b} \sim k^{a-b}
$$

we have

$$
\begin{gather*}
J_{1} \sim\left(\frac{k^{-1-m}}{m!}\right)^{2} \sum_{n=0}^{\infty} A_{n} \Gamma(4 \nu-2 m-1)(2 k)^{2 m-4 \nu-1} \sim k^{-4 \nu-1} 2^{2 m-4 \nu-1} \Gamma(4 \nu-2 m-1) \sum_{n=0}^{\infty} \frac{A_{n}}{m!} .  \tag{45}\\
J_{2}=\mathcal{O}_{k \sim \infty}(1) \tag{46}
\end{gather*}
$$

In the same

$$
\begin{equation*}
J_{3} \sim k^{m-4 \nu+1} \frac{\alpha_{k}^{\nu, m} \Gamma(4 \nu-2 m-1)}{m!(2 \nu-m-1)} \sum_{n=0}^{\infty} A_{n} . \tag{47}
\end{equation*}
$$

Therefore

$$
\lambda_{k} \sim C \sqrt{k^{m-4 \nu+1}}
$$

where $C$ is a constant.

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## Chapter 2

# Some Tauberian Theorems under Triple Statistically Nörlund-Cesáro Summability Method 

Carlos Granados


#### Abstract

In this paper, we extend the notion presented by Braha (2020) in a higher dimension, we introduce the notion of $N_{p, q}^{n, m, g} C_{n, m, g}^{(1,1,1)}$-statistically convergence and show necessity and sufficiency conditions under which the existence of the limit $s t-\lim _{n, m, g \rightarrow \infty} x_{n, m, g}=L$ follows from that $s t-\lim _{n, m, g \rightarrow \infty} N_{p, q}^{n, m, g} C_{n, m g}^{(1,1,1)}=L$. These conditions are one-sided or two-sided if $\left(x_{n, m, g}\right)$ is a sequence of real or complex numbers, respectively.


Keywords: Nörlund-Cesro summability method, one-sided and two-sided Tauberian conditions, triple statistical convergence

## 1. Introduction

The concept of statistical convergence was introduced by Fast [1] and Steinhaus [2]. Besides, in this connection, Fridy [3] showed some relation to a Tauberian condition for the statistical convergence of $\left(x_{k}\right)$. Subsequently, many researchers have worked in this area in several settings. For more recent works in this direction, one may refer to [4, 5]. Existing works in this field based on statistical convergence appears to have been restricted to real or complex sequences; however, Parida et al. [6] extended the idea for a locally convex Hausdorff topological linear space. Tauber [7] introduced the first Tauberian theorems for single sequences, that an Abel summable sequence is convergent with some suitable conditions. Later, a huge number of authors such as Landau [8], Hardy and Littlewood [9], and Schmidt [10] obtained some classical Tauberian theorems for Cessáro and Abel summability methods of single sequences. Recently, Braha [11] introduced some notions on statistical convergence by using the Nörlund-Cesáro summability method in a single sequence and proved some Tauberian theorems. In the last year, Canak and Totur [12], and Jena et al. [13] presented and studied several Tauberian theorems for single sequences. On the other hand, Knopp [14] obtained some classical type Tauberian theorems for Abel and $(C, 1,1)$ summability methods of double sequences and proved that Abel and $(C, 1,1)$ summability methods hold for the set of bounded sequences. Further, Moricz [15] proved some Tauberian theorems for Cesáro summable double sequences and deduced Tauberian theorems of Landau [16] and Hardy [17] type. Canak and Totur [18]
have proved a Tauberian theorem for Cesáro summability of single integrals and also the alternative proofs of some classical type Tauberian theorems for the Cesáro summability of single integrals and later introduced by Parida et al. [6] for double integrals. Otherwise, the notion of $(C, 1,1,1)$ summability of a triple sequence was originally introduced by Canak and Totur in 2016 [19]. Later, Canak et al. [20] studied some $(C, 1,1,1)$ means of a statistical convergent triple sequence and gave some classical Tauberian theorems for a triple sequence that $P$-convergence follows from statistically $(C, 1,1,1)$ summability under the two-sided boundedness conditions and slowly oscillating conditions in certain senses. Then, in 2020 Totur and Canak [21] proved Tauberian conditions under which convergence of triple integrals follows from ( $C, 1,1,1$ ) summability. For more studies associated to the main topic of this paper, we refer the reader to [22-24].

Let $\left(p_{n, m, g}\right)$ and $\left(q_{n, m, g}\right)$ be any two non-negative real sequences with

$$
R_{n, m, g}=\sum_{i=0}^{n} \sum_{j=0}^{m} \sum_{k=0}^{g} p_{i, j, k} q_{n-i, m-j, g-k} \neq 0((n, m, g) \in \mathbb{N} \times \mathbb{N} \times \mathbb{N})
$$

and ( $C, 1,1,1$ )-Cesáro summability method. Let $\left(x_{n, m, g}\right)$ be a sequence of real of complex numbers and set

$$
N_{p, q}^{n, m, g} C_{n, m, g}^{(1,1,1)}=\frac{1}{R_{n, m, g}} \sum_{i=0}^{n} \sum_{j=0}^{m} \sum_{k=0}^{g} p_{i, j, k} q_{n-i, m-j, g-k} \frac{1}{i+1} \frac{1}{j+1} \frac{1}{k+1} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k} x_{u, v, y}
$$

for $(n, m, g) \in \mathbb{N} \times \mathbb{N} \times \mathbb{N}$.
In this paper, we show necessary and sufficient conditions under which the existence of the limit $\lim _{n, m, g \rightarrow \infty} x_{n, m, g}=L$ follows from that of $\lim _{n, m, g \rightarrow \infty} N_{p, q}^{n, m, g} C_{n, m, g}^{(1,1,1)}=L$. These conditions are one-sided or two-sided if $\left(x_{n, m, g}\right)$ is a sequence of real or complex numbers, respectively.

Given two non-negative sequences $\left(p_{n, m, g}\right)$ and $\left(q_{n, m, g}\right)$, the convolution $(p \star q)$ is defined by

$$
R_{n, m, g}=(p \star q)_{n, m, g}=\sum_{i=0}^{n} \sum_{j=0}^{m} \sum_{k=0}^{g} p_{i, j, k} q_{n-i, m-j, g-k}=\sum_{i=0}^{n} \sum_{j=0}^{m} \sum_{k=0}^{g} p_{n-i, m-j, g-k} q_{i, j, k}
$$

with $(C, 1,1,1)$ we will denote the triple Cesáro summability method. Now, let $\left(x_{n, m, g}\right)$ be a sequence, when $(p \star q)_{n, m, g} \neq 0$ for all $(n, m, g) \in \mathbb{N} \times \mathbb{N} \times \mathbb{N}$ the generalized Nörlund-Cesáro transform of the sequence $\left(x_{n, m, g}\right)$ is the sequence $N_{p, q}^{n, m, g} C_{n, m, g}^{(1,1,1)}$ obtained by putting

$$
\begin{equation*}
N_{p, q}^{n, m, g} C_{n, m g}^{(1,1,1)}=\frac{1}{(p \star q)_{n, m, g}} \sum_{i=0}^{n} \sum_{j=0}^{m} \sum_{k=0}^{g} p_{i, j, k} q_{n-i, m-j, g-k} \frac{1}{i+1} \frac{1}{j+1} \frac{1}{k+1} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k} x_{u, v, y} . \tag{1}
\end{equation*}
$$

We say that the sequence $\left(x_{n, m, g}\right)$ is generalized Nörlund-Cesáro summable to $L$ determined by the sequences $\left(p_{n, m, g}\right)$ and $\left(q_{n, m, g}\right)$ (or simply summable $N_{p, q}^{n, m, g} C_{n, m, g}^{(1,1,1)}$ ) to $L$ if

$$
\begin{equation*}
\lim _{n, m, g \rightarrow \infty} N_{p, q}^{n, m, g} C_{n, m, g}^{(1,1,1)}=L . \tag{2}
\end{equation*}
$$

Throughout this paper, we will assume that the sequences $\left(p_{n, m, g}\right)$ and $\left(q_{n, m, g}\right)$ are satisfying the following conditions

$$
\begin{gather*}
q_{n, m, g} \geq 1, \quad \sum_{i=0}^{n} \sum_{j=0}^{m} \sum_{k=0}^{g} p_{i, j} \sim n m g, \quad(n, m, g) \in \mathbb{N} \times \mathbb{N} \times \mathbb{N},  \tag{3}\\
q_{\lambda_{n-i, m-j,-k}} \leq 2 q_{n-i, m-j, g-k}, i=1,2, \ldots, n ; j=1,2, \ldots, m ; \lambda>1 k=1,2, \ldots ; \lambda>1,  \tag{4}\\
q_{n-i, m-j, g-k} \leq 2 q_{\lambda_{n-i, m-j-k}} i=1,2, \ldots, \lambda_{n} ; j=1,2, \ldots, \lambda_{m} ; k=1,2, \ldots, ; 0<\lambda<1, \tag{5}
\end{gather*}
$$

where $\lambda_{n}=[\lambda n], \lambda_{m}=[\lambda m]$ and $\lambda_{g}=[\lambda g]$. On the other hand, $a_{n, m g} \sim b_{n, m, g}$ means that there are constants $C, C_{1}$ such that $a_{n, m, g} \leq C b_{n, m, g} \leq C_{1} a_{n, m, g}$. If

$$
\begin{equation*}
\lim _{n, m, g \rightarrow \infty} x_{n, m, g}=L \tag{6}
\end{equation*}
$$

implies (2), then the method $N_{p, q}^{n, m, g} C_{n, m, g}^{(1,1,1)}$ is said to be regular. Nevertheless, the converse is not always true as can be seen in the following example:

Let us consider that $p_{n, m, g}=q_{n, m g}=1$ for all $(n, m, g) \in \mathbb{N} \times \mathbb{N} \times \mathbb{N}$. Besides, we define the following sequence $x=\left(x_{i, j, k}\right)=(-1)^{i+j+k}$, then we get

$$
\begin{aligned}
& \frac{1}{(n+1)(m+1)(g+1)}\left|\sum_{i=0}^{n} \sum_{j=0}^{m} \sum_{k=0}^{g} \frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}(-1)^{u+v+y}\right| \\
& \leq \frac{1}{(n+1)(m+1)(g+1)} \sum_{i=0}^{n} \sum_{j=0}^{m} \sum_{k=0}^{g} \frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k} 1 \rightarrow 1 \text { as } n, m, g \rightarrow \infty .
\end{aligned}
$$

and as we know, $x=\left(x_{i, j, k}\right)$ is not convergent. Notice that (6) can imply (2) under a certain condition, which is called a Tauberian conditions. Any theorem which states that convergence of a sequence follows from its $N_{p, q}^{n, m, g} C_{n, m, g}^{(1,1,1)}$ summability and some Tauberian conditions are said to be a Tauberian theorems for the $N_{p, q}^{n, m, g} C_{n, m, g}^{(1,1,1)}$ summability method.

Next, we will find some conditions under which the converse implication holds, for defined convergence. Exactly, we will prove under which conditions statistical convergence of sequences $\left(x_{n, m, g}\right)$, follows from statistically Nörlund-Cesáro summability method.

A sequence $\left(x_{n, m g}\right)$ is weighted $N_{p, q}^{n, m, g} C_{n, m g}^{(1,1,1)}$-statistically convergent to $L$ if for every $\varepsilon>0$,

$$
\begin{aligned}
& \left.\lim _{n, m, g \rightarrow \infty} \frac{1}{(p \star q)_{n, m, g}} \right\rvert\,\left\{i, j, k \leq(p \star q)_{n, m, g}:\right. \frac{1}{(p \star q)_{n, m g}} \sum_{i=0}^{n} \sum_{j=0}^{m} \sum_{k=0}^{g} p_{i, j, k} q_{n-i, m-j, g-k} \\
&\left.\left.\frac{1}{i+1} \frac{1}{j+1} \frac{1}{(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k} x_{u, v, y}-L \right\rvert\, \geq \varepsilon\right\} \mid=0 .
\end{aligned}
$$

And we say that the sequence $\left(x_{n, m, g}\right)$ is statistically summable to $L$ by the weighted summability method $N_{p, q}^{n, m, g} C_{n, m g}^{(1,1,1)}$ if $s t-\lim _{n, m, g} N_{p, q}^{n, m, g} C_{n, m, g}^{(1,1,1)}=L$. We will denote by $N_{p, q}^{n, m, g} C_{n, m g}^{(1,1,1)}(s t)$ the set of all sequences which are statistically summable, by the weighted summability method $N_{p, q}^{n, m, g} C_{n, m g}^{(1,1,1)}$.

Theorem 1.1 Let $x=\left(x_{n, m, g}\right)$ be a sequence $N_{p, q}^{n, m, g} C_{n, m, g}^{(1,1,1)}$ summable to $L$, then the sequence $x=\left(x_{n, m, g}\right)$ is $N_{p, q}^{n, m, g} C_{n, m, g}^{(1,1,1)}$-statistically convergent to $L$, but not conversely.

Proof: The first part of the proof is obvious. To prove the second part, we will show the following example:

Let us define

$$
x_{i, j, k}=\left\{\begin{array}{lc}
\sqrt{x y z}, & \text { for } i=n^{2} j=m^{2} \text { and } k=g^{2} \\
0, & \text { otherwise }
\end{array}\right.
$$

and $p_{n, m, g}=1=q_{n, m, g}$. Under this conditions we obtain,

$$
\begin{aligned}
& \left.\frac{1}{(n+1)(m+1)(g+1)} \right\rvert\,\{i, j, k \leq n+1, m+1, g+1: \\
& \left.\left|\frac{1}{(n+1)(m+1)(g+1)} \sum_{i=0}^{n} \sum_{j=0}^{m} \sum_{k=0}^{g} \frac{1}{P_{i, j, k}} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k} p_{u, v, y} x_{u, v, y}-0\right| \geq \varepsilon\right\} \mid \\
& \leq \frac{\sqrt{(n+1)(m+1)(g+1)}}{(n+1)(m+1)(g+1)} \rightarrow 0 .
\end{aligned}
$$

On the other hand, for $i=n^{2}, j=m^{2}$ and $k=g 2$, we have

$$
\begin{aligned}
& \frac{1}{(n+1)(m+1)(g+1)} \sum_{i=0}^{n} \sum_{j=0}^{m} \sum_{k=0}^{g} \frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k} x_{u, v, y} \rightarrow \infty, \\
& \text { as } n, m, g \rightarrow \infty .
\end{aligned}
$$

From last relation follows that $x=\left(x_{n, m, g}\right)$ is not $N_{p, q}^{n, m, g} C_{n, m, g}^{(1,1,1)}$ summable to 0 .
Theorem 1.2 Let $x=\left(x_{n, m g}\right)$ be a sequence statistically convergent to $L$ and $\left|x_{n, m, g}-L\right| \leq M$ for every $(n, m, g) \in \mathbb{N} \times \mathbb{N} \times \mathbb{N}$. Then, it converges $N_{p, q}^{n, m, g} C_{n, m, g}^{(1,1,1)}$-statistically to $L$.

Proof: From the fact that $\left(x_{n, m, g}\right)$ converges statistically to $L$, we have

$$
\lim _{n, m, g \rightarrow \infty} \frac{\left|i, j, k \leq n, m, g:\left|x_{i, j, k}-L\right| \geq \varepsilon\right\} \mid}{n m g}=0 .
$$

We will denote $B_{\varepsilon}=\left\{i, j, k \leq n, m, g:\left|x_{i, j, k}-L\right| \geq \varepsilon\right\}$ and $\bar{B}_{\varepsilon}=$ $\left\{i, j, k \leq n, m, g:\left|x_{i, j, k}-L\right| \leq \varepsilon\right\}$. Then,

$$
\begin{aligned}
& \left|\frac{1}{R_{n, m, g}} \sum_{i=0}^{n} \sum_{j=0}^{m} \sum_{k=0}^{g} p_{i, j, k} q_{n-i, m-j, g-k} \frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k} x_{u, v, y}-L\right| \\
& =\left|\frac{1}{R_{n, m, g}} \sum_{i=0}^{n} \sum_{j=0}^{m} \sum_{k=0}^{g} p_{i, j, k} q_{n-i, m-j, g-k} \frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left(x_{u, v, y}-L\right)\right|
\end{aligned}
$$

$$
\begin{aligned}
& \leq \frac{1}{R_{n, m g}} \sum_{i=0}^{n} \sum_{j=0}^{m} \sum_{k=0}^{g} \quad p_{i, j, k} q_{n-i, m-j, g-k} \frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left|x_{u, v, y}-L\right| \\
& +\frac{1}{R_{n, m, g}} \sum_{i=0}^{n} \sum_{j=0}^{m} \sum_{k=0}^{g} p_{i, j, k} q_{n-i, m-j, g-k} \frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left|x_{u, v, y}-L\right| \\
& \leq M \frac{1}{R_{n, m, g}} \sum_{i=0}^{n} \sum_{j=0}^{m} \sum_{k=0}^{g} \quad 1+\varepsilon \leq M \frac{C_{2}}{n m g} \sum_{i=0}^{n} \sum_{j=0}^{m} \sum_{k=0}^{g} \quad 1+\varepsilon \rightarrow 0+\varepsilon,
\end{aligned}
$$

$$
\text { as } n, m, g \rightarrow \infty \text {, }
$$

for some constant $C_{2}$.
Converse of Theorem 1.2 is not true as can be seen in the following example.
Consider that $p_{n, m, g}=(n+1)(m+1)(g+1),\left(q_{n, m, g}\right)=1$ for some $(n, m, g) \in \mathbb{N} \cup\{0\} \times \mathbb{N} \cup\{0\} \times \mathbb{N} \cup\{0\}$ and define the sequence $x=\left(x_{n, m, g}\right)$ as follows:
$x_{i, j, k}=\left\{\begin{array}{l}1, \quad \text { for } i=p^{2}-p, \ldots, p^{2}-j=t^{2}-t, \ldots, t^{2}-1 \text { and } k=o^{2}-o, \ldots, o^{2}-1 ; \\ -\frac{1}{p t o}, \quad \text { for } i=p^{2}, p=2, . . j=t^{2}, t=2, . . \text { and } k=o^{2}, o=2, \ldots \\ 0, \\ \text { otherwise }\end{array}\right.$
Under this conditions, after some basic calculations we get that $x=\left(x_{n, m g}\right)$ is $N_{p, q}^{n, m, g} C_{n, m g}^{(1,1,1)}$-summable to 1 . Therefore, by Theorem 1.2, $x=\left(x_{n, m, g}\right)$ is $N_{p, q}^{n, m, g} C_{n, m, g}^{(1,1,1)}-$ statistically convergent. On the other hand, the sequences $p^{2} ; p=2,3, \ldots, t^{2} ; t=$ $2,3, \ldots$ and $o^{2} ; o=2,3, \ldots$ have natural density zero and it is clear that st$\lim _{n, m, g} \inf _{n, m, g}=0$ and $s t-\lim _{n, m, g} \sup _{n, m, g}=1$. Hence, $\left(x_{i, j, k}\right)$ is not statistically convergent.

## 2. Tauberian theorems under $N_{p, q}^{n, m, g} C_{n, m, g}^{(1,1,1)}$-statistically convergence

In this section, we show the results that we obtained. Throughout this paper, $R_{\lambda_{n, m g}}$ and $R_{\lambda_{n}, \lambda_{m}, \lambda_{g}}$ will have the same meaning.

Consider that $s t$ - $\lim i, j, k x_{i, j, k}=L ;\left(x_{n, m, g}\right)$ is $N_{p, q}^{n, m, g} C_{n, m, g}^{(1,1,1)}$-statistically convergent and (13) satisfies, then for every $t>1$, is valid the following relation

$$
\begin{align*}
& s t-\lim _{i, j, k} \frac{1}{R_{\lambda_{i, j, k}}-R_{i, j, k}} \sum_{w=i+1}^{\lambda_{i}} \sum_{e=j+1}^{\lambda_{j}} \sum_{r=k+1}^{\lambda_{k}} p_{w, e, r} q_{\lambda_{i}-w, \lambda j-e, \lambda k-r}  \tag{7}\\
& \frac{1}{(w+1)(e+1)(r+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left(x_{w, e, r}-x_{i, j, k}\right)=0
\end{align*}
$$

and in case where $0<t<1$,

$$
\begin{align*}
& s t-\lim _{i, j, k} \frac{1}{R_{i, j, k}-R_{\lambda_{i, j, k}}} \sum_{w=\lambda_{i}+1}^{i} \sum_{e=\lambda_{j}+1}^{j} \sum_{r=\lambda_{k}+1}^{k} p_{w, e, r} q_{i-w, j-e, k-r} \\
& \frac{1}{(w+1)(e+1)(r+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left(x_{i, j, k}-x_{w, e, r}\right)=0 . \tag{8}
\end{align*}
$$

The condition given by relation (13) is equivalent to the condition

$$
\begin{equation*}
s t-\lim _{n, m, g \rightarrow \infty} \frac{R_{n, m, g}}{R_{n, m, g}}>1, \quad 0<\lambda<1 \tag{9}
\end{equation*}
$$

Proof: Suppose that relation (13) is valid, $0<\lambda<1, w=\lambda_{n}=[\lambda n], e=\lambda_{m}=[\lambda m]$ and $r=\lambda_{g}=[\lambda g],(n, m, g) \in \mathbb{N} \times \mathbb{N} \times \mathbb{N}$. Then, it follows that

$$
\frac{1}{\lambda}>1 \Rightarrow \frac{w}{\lambda}=\frac{[\lambda n]}{t} \leq n, \frac{1}{\lambda}>1 \Rightarrow \frac{e}{\lambda}=\frac{[\lambda m]}{t} \leq m \text { and } \frac{1}{\lambda}>1 \Rightarrow \frac{r}{\lambda}=\frac{[\lambda g]}{t} \leq g .
$$

From above relation and definition of sequences $\left(p_{n, m, g}\right)$ and $\left(q_{n, m, g}\right)$, we have

Conversely, suppose that (9) is valid. Now, let $\lambda>1$ be given and let $\lambda_{1}, \lambda_{2}, \lambda_{3}$ be chosen such that $1<\lambda_{1}, \lambda_{2}, \lambda_{3}<\lambda$. Set $w=\lambda_{n}=[\lambda n], e=\lambda_{m}=[\lambda m]$ and $r=\lambda_{g}=[\lambda g]$. From $0<\frac{1}{\lambda}<\frac{1}{\lambda_{1}}, \frac{1}{\lambda_{2}}, \frac{1}{\lambda_{3}}<1$, it follows that

$$
n \leq \frac{\lambda n-1}{\lambda_{1}}<\frac{[\lambda n]}{\lambda_{1}}=\frac{w}{\lambda_{1}}, m \leq \frac{\lambda m-1}{\lambda_{2}}<\frac{[\lambda m]}{\lambda_{2}}=\frac{e}{\lambda_{2}} \text { and } g \leq \frac{\lambda g-1}{\lambda_{3}}<\frac{[\lambda g]}{\lambda_{3}}=\frac{r}{\lambda_{3}}
$$

provided $\lambda_{1}, \lambda_{2}, \lambda_{3} \leq \lambda-\frac{1}{n}, \lambda-\frac{1}{m}, \lambda-\frac{1}{g}$, which is a case where if $n, m$ and $g$ are large enough. Under this condition, we obtain

$$
\frac{R_{\lambda_{n, m g}}}{R_{n, m, g}} \geq \frac{R_{\lambda_{n, m g}}}{R_{\left[\frac{w}{\lambda_{1}}\right],\left[\frac{e}{\lambda_{2}}\right],\left[\frac{r}{k_{3}}\right]}} \Rightarrow s t-\operatorname{liminn}_{n, m, g \rightarrow \infty} \frac{R_{\lambda_{n, m g}}}{R_{n, m, g}} \geq s t-\liminf _{n, m, g \rightarrow \infty} \frac{R_{\lambda_{n, m g}}}{R_{\left[\frac{w}{1_{1}}\right]}\left[\frac{e}{\left[\frac{e}{2}\right.}\right],\left[\frac{r}{\left[\frac{r}{3}\right.}\right]}>1
$$

Consider that (13) is satisfied and let $x=\left(x_{i, j, k}\right)$ be a sequence of complex numbers which is $N_{p, q}^{n, m g} C_{n, m, g}^{(1,1)}$-statistically convergent to $L$. Then,

$$
\begin{align*}
& s t-\lim _{n, m, g} \frac{1}{R_{\lambda_{n, m g}}-R_{n, m, g}} \sum_{i=n+1}^{\lambda_{n}} \sum_{j=m+1}^{\lambda_{m}} \sum_{k=g+1}^{\lambda_{g}} p_{i, j, k} q_{\lambda_{n}-i, \lambda_{m}-j, \lambda_{g}-k}  \tag{10}\\
& \frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k} x_{u, v, y}=L \text { for } \lambda>1
\end{align*}
$$

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and

$$
\begin{align*}
& s t-\lim _{n, m, g} \frac{1}{R_{n, m, g}-R_{\lambda_{n, m g}}} \sum_{i=\lambda_{n}+1}^{n} \sum_{j=\lambda_{m}+1}^{m} \sum_{k=\lambda_{g}+1}^{g} p_{i, j, k} q_{n-i, m-j, g-k}  \tag{11}\\
& \frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k} x_{u, v, y}=L \text { for } 0<\lambda<1 .
\end{align*}
$$

Proof: We begin proving the case (10), i.e. when $\lambda>1$. Then, we have

$$
\begin{align*}
& \frac{1}{R_{\lambda_{n, m g}}-R_{n, m, g}} \sum_{i=n+1}^{\lambda_{n}} \sum_{j=m+1}^{\lambda_{m}} \sum_{k=g+1}^{\lambda_{g}} p_{i, j, k} q_{\lambda_{n}-i, \lambda_{m}-j, \lambda_{g}-k} \frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left(x_{u, v, y}-L\right) \\
& =\frac{R_{\lambda_{n, m g}}}{R_{\lambda_{n, m g}}-R_{n, m, g}} \frac{1}{R_{\lambda_{n, m g}}} \sum_{i=n+1}^{\lambda_{n}} \sum_{j=m+1}^{\lambda_{m}} \sum_{k=g+1}^{\lambda_{g}} p_{i, j, k} q_{\lambda_{n}-i, \lambda_{m}-j, \lambda_{g}-k} \frac{1}{(i+1)(j+1)(k+1)} \\
& \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left(x_{u, v, y}-L\right)-\frac{R_{n, m, g}}{R_{\lambda_{n, m, g}}-R_{n, m, g}} \frac{1}{R_{n, m g}} \sum_{i=0}^{n} \sum_{j=0}^{m} \sum_{k=0}^{g} p_{i, j, k} q_{\lambda_{n}-i, \lambda_{m}-j, \lambda_{g}-k} \\
& \frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left(x_{u, v, y}-L\right) \\
& =\frac{R_{\lambda_{n, m, g}}}{R_{\lambda_{n, m g}}-R_{n, m, g}} \frac{1}{R_{\lambda_{n, m g}, m}} \sum_{i=n+1}^{\lambda_{n}} \sum_{j=m+1}^{\lambda_{m}} \sum_{k=g+1}^{\lambda_{g}} p_{i, j, k} q_{\lambda_{n}-i, \lambda_{m}-j, \lambda_{g}-j} \\
& \frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left(x_{u, v, y}-L\right)-\frac{R_{n, m, g}}{R_{\lambda_{n, m, g}-R_{n, m, g}}} \frac{1}{R_{n, m, g}} \sum_{i=0}^{n} \sum_{j=0}^{m} \sum_{k=0}^{g} p_{i, j, k} \\
& \left(q_{\lambda_{n}-i, \lambda_{m}-j, \lambda_{g}-k}+q_{n-i, m-j, g-k}-q_{n-i, m-j, g-k}\right) \frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left(x_{u, v, y}-L\right) \\
& =\frac{R_{\lambda_{n, m g}}}{R_{\lambda_{n, m g}}-R_{n, m g}} \frac{1}{R_{\lambda_{n, m g}}} \sum_{i=n+1}^{\lambda_{n}} \sum_{j=m+1}^{\lambda_{m}} \sum_{k=g+1}^{\lambda_{g}} p_{i, j, k} q_{\lambda_{n}-i, \lambda_{m}-j, \lambda_{g}-k} \\
& \frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left(x_{u, v, y}-L\right)-\frac{R_{n, m, g}}{R_{\lambda_{n, m, g}}-R_{n, m, g}} \frac{1}{R_{n, m, g}} \sum_{i=0}^{n} \sum_{j=0}^{m} \sum_{k=0}^{g} p_{i, j, k} \\
& q_{n-i, m-j, g-k} \frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k} x_{u, v, y}-\frac{R_{n, m, g}}{R_{\lambda_{n, m, g}}-R_{n, m, g}} \frac{1}{R_{n, m, g}} \sum_{i=0}^{n} \sum_{j=0}^{m} \sum_{k=0}^{g} p_{i, j, k} \\
& \left(q_{\lambda_{n}-i, \lambda_{m}-j, \lambda_{g}-k}-q_{n-i, m-j, g-k}\right) \frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left(x_{u, v, y}-L\right) . \tag{12}
\end{align*}
$$

From (12), definition of the sequence $\left(q_{n, m, g}\right)$ and relation $\limsup _{n, m, g} \frac{R_{2, m, m_{g}}}{R_{2 n, m, g}-R_{n, m, g}}<\infty$, we get (10).

Prove of (11) is made similarly to the prove of (10).
In the following theorem, we characterize the converse implication when the statistically convergence follows from its $N_{p, q}^{n, m, g} C_{n, m g}^{(1,1,1)} \hat{a}^{\wedge ’}$ statistically convergence.

Theorem 1.3 Let $\left(p_{n, m, g}\right)$ and $\left(q_{n, m, g}\right)$ be two non-negative real sequences and

$$
\begin{equation*}
s t-\liminf _{n, m, g \rightarrow \infty} \frac{R_{n, m g}}{R_{n, m, g}}>1 \text { for every } \lambda>1, \tag{13}
\end{equation*}
$$

where $\lambda_{n, m, g}=\lambda_{n} \lambda_{m} \lambda_{g}=[\lambda n][\lambda m][\lambda g]$ denotes the integral part of $\lambda n \lambda m \lambda g$ for every $(n, m, g) \in \mathbb{N} \times \mathbb{N} \times \mathbb{N}$, and let $\left(x_{n, m, g}\right)$ be a sequence of real numbers which is $N_{p, q}^{n, m, g} C_{n, m, g}^{(1,1,1)}$-statistically convergent to a finite number $L$. Then, $\left(x_{n, m, g}\right)$ is $s t$-convergent to the same number $L$ if and only if the following two conditions hold

$$
\begin{align*}
& \inf _{\lambda>1} \limsup _{n, m, g} \frac{1}{R_{n, m, g}} \left\lvert\,\left\{i, j, k \leq R_{n, m, g}: \frac{1}{R_{\lambda_{i, j, k}}-R_{i, j, k}} \sum_{w=i+1}^{\lambda_{i}} \sum_{e=j+1}^{\lambda_{j}} \sum_{r=k+1}^{\lambda_{k}} p_{w, e, r} q_{\lambda_{i}-w, \lambda_{j}-e, \lambda_{k}-r}\right.\right. \\
& \left.\frac{1}{(w+1)(e+1)(r+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left(x_{w, e, r}-x_{i, j, k}\right) \leq-\varepsilon\right\} \mid=0, \tag{14}
\end{align*}
$$

and

$$
\begin{align*}
& \inf _{0<\lambda<1} \limsup _{n, m, g} \frac{1}{R_{n, m, g}} \left\lvert\,\left\{i, j, k \leq R_{n, m, g}: \frac{1}{R_{i, j, k}-R_{\lambda_{i, j, k}}} \sum_{w=\lambda_{i}+1}^{i} \sum_{e=\lambda_{j}+1}^{j} \sum_{r=\lambda_{k}+1}^{k} p_{w, e, r} q_{i-w, j-e, k-r}\right.\right. \\
& \left.\frac{1}{(w+1)(e+1)(r+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left(x_{i, j, k}-x_{w, e, r}\right) \leq-\varepsilon\right\} \mid=0 . \tag{15}
\end{align*}
$$

Proof: Necessity: Suppose that $\lim _{n, m g \rightarrow \infty} x_{n, m, g}=L$ and (13) holds. By Proposition 2, we have

$$
\begin{aligned}
& \lim _{n, m, g \rightarrow \infty} \frac{1}{R_{\lambda_{n, m g}}-R_{n, m, g}} \sum_{i=n+1}^{\lambda_{n}} \sum_{j=m+1}^{\lambda_{m}} \sum_{k=g+1}^{\lambda_{g}} p_{i, j, k} q_{\lambda_{n}-i, \lambda_{m}-j, \lambda_{g}-k} \\
& \frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left(x_{u, v, y}-x_{n, m, g}\right) \\
& =\lim _{n, m, g \rightarrow \infty}\left\{\left(\frac{1}{R_{\lambda_{n, m g}}-R_{n, m, g}} \sum_{i=n+1}^{\lambda_{n}} \sum_{j=m+1}^{\lambda_{m}} \sum_{k=g+1}^{\lambda_{g}} p_{i, j, k} q_{\lambda_{n}-i, \lambda_{m}-j, \lambda_{g}-k}\right.\right. \\
& \left.\left.\frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k} x_{u, v, y}\right)-x_{n, m, g}\right\}=0,
\end{aligned}
$$

for every $\lambda>1$. In case where $0<\lambda<1$, we have that

$$
\begin{aligned}
& \lim _{n, m, g \rightarrow \infty} \frac{1}{R_{n, m, g}-R_{\lambda_{n, m g}}} \sum_{i=\lambda_{n}+1}^{n} \sum_{j=\lambda_{m}+1}^{m} \sum_{k=\lambda_{g}+1}^{g} p_{i, j, k} q_{n-i, m-j, g-k} \\
& \frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left(x_{n, m, g}-x_{u, v, y}\right) \\
& =\lim _{n, m, g \rightarrow \infty}\left\{x_{n, m, g}-\left(\frac{1}{R_{n, m, g}-R_{\lambda_{n, m g}}} \sum_{i=\lambda_{n}+1}^{n} \sum_{j=\lambda_{m}+1}^{m} \sum_{k=\lambda_{g}+1}^{g} p_{i, j, k} q_{n-i, m-j, g-k}\right.\right. \\
& \left.\left.\frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k} x_{u, v, y}\right)\right\}=0 .
\end{aligned}
$$

Sufficiency: Consider that (14) and ((15) are satisfied. In what follows, we will prove that $\lim _{n, m, g \rightarrow \infty} x_{n, m, g}=L$. Given any $\varepsilon>0$, by (14) we can choose $\lambda_{1}>0$ such that

$$
\begin{align*}
& \liminf _{n, m, g \rightarrow \infty} \frac{1}{R_{\lambda_{1}}, \lambda_{m_{1}}, \lambda_{k_{1}}-R_{n, m g}} \sum_{i=n+1}^{\lambda_{n_{1}}} \sum_{j=m+1}^{\lambda_{m_{1}}} \sum_{k=g+1}^{\lambda_{g_{1}}} p_{i, j, k} q_{\lambda_{n}-i, \lambda_{m}-j, \lambda_{g}-k}  \tag{16}\\
& \frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left(x_{u, v, y}-x_{n, m, g}\right) \geq-\varepsilon,
\end{align*}
$$

where $\lambda_{n_{1}}=\left[\lambda n_{1}\right], \lambda_{m_{1}}=\left[\lambda m_{1}\right]$ and $\lambda_{g_{1}}=\left[\lambda g_{1}\right]$. By the assumed summability $N_{p, q}^{n, m, g} C_{n, m, g}^{(1,1,1)}$ of $\left(x_{n, m, g}\right)$, Proposition 2 and (16), we have

$$
\begin{equation*}
\limsup _{n, m, g \rightarrow \infty} x_{n, m, g} \leq L+\varepsilon, \tag{17}
\end{equation*}
$$

for any $\lambda>1$.
On the other hand, if $0<\lambda<1$, for every $\varepsilon>0$, we can choose $0<\lambda_{2}<1$ such that

$$
\begin{align*}
& m \liminf _{n, m, g \rightarrow \infty} \frac{1}{R_{n, m, g}-R_{\lambda_{n_{2}}, \lambda_{m_{2}}, \lambda_{g 2}}} \sum_{i=\lambda_{n_{2}}+1}^{n} \sum_{j=\lambda_{m_{2}}+1}^{m} \sum_{k=\lambda_{g_{2}}+1}^{g} p_{i, j, k} q_{n-i, m-j, g-k}  \tag{18}\\
& \frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left(x_{n, m, g}-x_{u, v, y}\right) \geq-\varepsilon,
\end{align*}
$$

where $\lambda_{n_{2}}=\left[\lambda n_{2}\right], \lambda_{m_{2}}=\left[\lambda m_{2}\right]$ and $\lambda_{g_{2}}=\left[\lambda g_{2}\right]$. By the assumed summability $N_{p, q}^{n, m, g} C_{n, m, g}^{(1,1,1)}$ of $\left(x_{n, m, g}\right)$, Proposition 2 and (18), we have

$$
\begin{equation*}
\liminf _{n, m, g \rightarrow \infty} x_{n, m, g} \geq L-\varepsilon, \tag{19}
\end{equation*}
$$

for any $0<\lambda<1$.
Since $\varepsilon>0$ is arbitrary, combining (17) and (19), we obtain

$$
\lim _{n, m, g \rightarrow \infty} x_{n, m, g}=L .
$$

In the following theorem, we will consider the case where $x=\left(x_{n, m, g}\right)$ is a sequence of complex numbers.

Theorem 1.4 Let (13) be satisfied and let ( $x_{n, m g}$ ) be a sequence of complex numbers which is $N_{p, q}^{n, m, g} C_{n, m g}^{(1,1,1)}$-statistically convergent to a finite number $L$. Then, $\left(x_{n, m, g}\right)$ is convergent to the same number $L$ if and only if the following two conditions hold

$$
\begin{align*}
& \inf _{\lambda>1} \limsup _{n, m, g} \frac{1}{R_{n, m, g}} \left\lvert\,\left\{i, j, k \leq R_{n, m, g}: \frac{1}{R_{\lambda_{i, j, k}}-R_{i, j, k}} \sum_{w=i+1}^{\lambda_{i}} \sum_{e=j+1}^{\lambda_{j}} \sum_{r=k+1}^{\lambda_{k}} p_{w, e, r} q_{\lambda_{i}-w, \lambda_{j}-e, \lambda_{k}-r}\right.\right. \\
& \left.\frac{1}{(w+1)(e+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left(x_{w, e, r}-x_{i, j, k}\right) \geq \varepsilon\right\} \mid=0, \tag{20}
\end{align*}
$$

and

$$
\begin{align*}
& \inf _{0<\lambda<1} \limsup _{n, m, g} \frac{1}{R_{n, m, g}} \left\lvert\,\left\{i, j, k \leq R_{n, m, g}: \frac{1}{R_{i, j, k}-R_{\lambda_{i, j, k}}} \sum_{w=\lambda_{i}+1}^{i} \sum_{e=\lambda_{j}+1}^{j} \sum_{r=\lambda_{k}+1}^{k} p_{w, e, r} q_{i-w, j-e, k-r}\right.\right. \\
& \left.\frac{1}{(w+1)(e+1)(r+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left(x_{i, j, k}-x_{w, e, r}\right) \geq \varepsilon\right\} \mid=0 . \tag{21}
\end{align*}
$$

Proof: Necessity: If both (2) and (6) hold, then Proposition 2 yields (20) for every $\lambda>1$ and (21) for every $0<\lambda<1$.

Sufficiency: Suppose that (2), (13) and one of the conditions (20) and (21) are satisfied. For any given $\varepsilon>0$, there exists $\lambda_{1}>0$ such that

$$
\begin{aligned}
& \limsup _{n, m, g \rightarrow \infty} \left\lvert\, \frac{1}{R_{\lambda_{n_{1}}, \lambda_{m_{1}}}, \lambda_{g_{1}}}-R_{n, m}\right. \\
& \sum_{i=n+1}^{\lambda_{n_{1}}} \sum_{j=m+1}^{\lambda_{m_{1}}} \sum_{k=g+1}^{\lambda_{g_{1}}} p_{i, j, k} q_{\lambda_{n_{1}}-i, \lambda_{m_{1}}-j, \lambda_{g_{1}}-k} \\
& \left.\frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left(x_{u, v, y}-x_{n, m, g}\right) \right\rvert\, \leq \varepsilon,
\end{aligned}
$$

where $\lambda_{n_{1}}=\left[\lambda n_{1}\right], \lambda_{m_{1}}=\left[\lambda m_{1}\right]$ and $\lambda_{g_{1}}=\left[\lambda g_{1}\right]$. Taking into account fact that $\left(x_{n, m, g}\right)$ is $N_{p, q}^{n, m, g} C_{n, m, g}^{(1,1,1)}$ summbale to $L$ and Proposition 2, we have the following estimation

$$
\begin{aligned}
& \limsup _{n, m, g \rightarrow \infty}\left|L-x_{n, m g}\right| \\
& \leq \limsup _{n, m, g \rightarrow \infty} \left\lvert\, L-\frac{1}{R_{\lambda_{n_{1}}, \lambda_{m_{1}}, \lambda_{g_{1}}}-R_{n, m}} \sum_{i=n+1}^{\lambda_{n_{1}}} \sum_{j=m+1}^{\lambda_{m_{1}}} \sum_{k=g+1}^{\lambda_{g_{1}}} p_{i, j, k} q_{\lambda_{n_{1}}-i, \lambda_{m_{1}}-j, \lambda_{g_{1}}-k}\right. \\
& \left.\frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k} x_{u, v, y} \right\rvert\, \\
& +\limsup _{n, m, g \rightarrow \infty} \left\lvert\, L-\frac{1}{R_{\lambda_{n_{1}}, \lambda_{m_{1}}, \lambda_{g_{1}}}-R_{n, m}} \sum_{i=n+1}^{\lambda_{n_{1}}} \sum_{j=m+1}^{\lambda_{m_{1}}} \sum_{k=g+1}^{\lambda_{g_{1}}} p_{i, j, k} q_{\lambda_{n_{1}}-i, \lambda_{m_{1}}-j, \lambda_{g_{1}}-k}\right. \\
& \left.\frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left(x_{u, v, y}-x_{n, m, g}\right) \right\rvert\, \\
& \leq \varepsilon .
\end{aligned}
$$

For a given $\varepsilon>0$, there exists $\lambda_{2}>0$ such that

$$
\begin{aligned}
& \limsup _{n, m, g \rightarrow \infty} \left\lvert\, \frac{1}{R_{n, m, g}-R_{\lambda_{n_{2}}, \lambda_{m_{2}}, \lambda_{g_{2}}}} \sum_{i=\lambda_{n_{2}}+1}^{n} \sum_{j=\lambda_{m_{2}}+1}^{m} \sum_{k=\lambda_{g_{2}}+1}^{g} p_{i, j, k} q_{n-i, m-j, g-k}\right. \\
& \left.\frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left(x_{n, m, g}-x_{u, v, y}\right) \right\rvert\, \leq \varepsilon,
\end{aligned}
$$

where $\lambda_{n_{2}}=\left[\lambda n_{2}\right], \lambda_{m_{2}}=\left[\lambda m_{2}\right]$ and $\lambda_{g_{2}}=\left[\lambda g_{2}\right]$. Taking into account fact that $\left(x_{n, m, g}\right)$ is $N_{p, q}^{n, m, g} C_{n, m, g}^{(1,1,1)}$ summbale to $L$ and Proposition 2, we obtain the following

$$
\begin{aligned}
& \limsup _{n, m, g \rightarrow \infty}\left|L-x_{n, m, g}\right| \\
& \limsup _{n, m \rightarrow \infty} \left\lvert\, L-\frac{1}{R_{n, m, g}-R_{\lambda_{n_{2}}, \lambda_{m_{2}}, \lambda_{g 2}}} \sum_{i=\lambda_{n_{2}}+1}^{n} \sum_{j=\lambda_{m_{2}}+1}^{m} \sum_{k=\lambda_{g_{2}}+1}^{g} p_{i, j, k} q_{n-i, m-j, g-k}\right. \\
& \left.\frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k} x_{u, v, y} \right\rvert\, \\
& +\limsup _{n, m, g \rightarrow \infty} \left\lvert\, \frac{1}{R_{n, m, g}-R_{\lambda_{n_{2}}, \lambda_{m_{2}}, \lambda_{g 2}}} \sum_{i=\lambda_{n_{2}}+1}^{n} \sum_{j=\lambda_{m_{2}}+1}^{m} \sum_{k=\lambda_{g_{2}}+1}^{g} p_{i, j, k} q_{n-i, m-j, g-k}\right. \\
& \left.\frac{1}{(i+1)(j+1)(k+1)} \sum_{u=0}^{i} \sum_{v=0}^{j} \sum_{y=0}^{k}\left(x_{n, m g}-x_{u, v, y}\right) \right\rvert\, \\
& \leq \varepsilon .
\end{aligned}
$$

Since $\varepsilon>0$ in either case, we get

$$
\lim _{n, m, g \rightarrow \infty} x_{n, m, g}=L
$$

## 3. Conclusion

In this paper, we have defined and proved new Tauberian theorems under triple statistically Nörlund-Cesáro summability, as a consequence of results showed in 2, some theorems, lemmas and corollaries can be defined and proved similarly by using $(1,0,0),(0,1,0)$, and $(0,0,1)$ method of summability. It is well know that Tauberian theorems for single sequences of single variable have been achieved a high degree of development; however, it is still in its infancy for triple sequences. For that reason, the results established in this paper can be extended and studied in some inclusion, Tauberian type theorems and Tauberian convexity type for certain families of generalized Nörlund.

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## Chapter 3

# A Brief Look at the Calderón and Hilbert Operators 

Guillermo J. Flores


#### Abstract

The Calderón operator is the sum of the Hardy averaging operator and its adjoint, and plays an important role in the theory of real interpolation. On the other hand, the Hilbert operator arises from the continuous version of Hilbert's inequality. Both operators appear in different contexts and have numerous applications within harmonic analysis. In this chapter we will briefly review the Calderón and Hilbert operators, showing some of the most relevant results within functional analysis and finally we will present recent results on these operators within Fourier analysis.


Keywords: Calderón operator, Hilbert operator, Lebesgue spaces, Lipschitz spaces, BMO spaces, weighted inequalities, Calderón weights

## 1. Introduction

The Calderón and Hilbert operators are among the most relevant operators in harmonic analysis, arising from Hilbert's double series theorem which is one of the simplest and most beautiful in the theory of double series of positive terms. It was proved by Hilbert, in the course of his investigations in the theory of integral equations, that the series $\sum_{m, n \in \mathbb{N}} \frac{a_{m} a_{n}}{a_{m}+a_{n}}$, where $a_{n} \geq 0$ for all $n \in \mathbb{N}$, is convergent whenever $\sum_{n \in \mathbb{N}} a_{n}^{2}$ is convergent.

Other proofs of Hilbert's double series theorem and generalizations in different directions were studied and published over time by influential mathematicians such as H. Weyl, F. Wiener, J. Schur, Fejér and F. Riesz, Pólya and Szegö, Francis and Littlewood, G.H. Hardy and M. Riesz, among others.

In [1, 2], G.H. Hardy observed that Hilbert's theorem stated above is an immediate corollary of another theorem which has interest in itself. This theorem is as follows: If $a_{n} \geq 0$ for all $n \in \mathbb{N}$ and $\sum_{n \in \mathbb{N}} a_{n}^{2}$ is convergent, then $\sum_{n \in \mathbb{N}}\left(\frac{1}{n} \sum_{j=1}^{n} a_{j}\right)^{2}$ is also convergent.

The first extension of the just stated Hilbert's and Hardy's results in which we are interested is the following (see [3]): Let $1<p<\infty$ and $p^{\prime}=p /(p-1)$ (i.e. $p^{\prime}$ is the conjugate of $p$ ). If $\sum_{n=1}^{\infty} a_{n}^{p}$ and $\sum_{n=1}^{\infty} b_{n}^{p^{\prime}}$ are convergent, where $a_{n}$ and $b_{n}$ are nonnegative numbers for all $n \in \mathbb{N}$, then
$\sum_{m=1}^{\infty} \sum_{n=1}^{\infty} \frac{a_{m} b_{n}}{m+n} \leq \frac{\pi}{\sin (\pi / p)}\left(\sum_{m=1}^{\infty} a_{m}^{p}\right)^{1 / p}\left(\sum_{n=1}^{\infty} b_{n}^{p^{\prime}}\right)^{1 / p^{\prime}}$ and $\sum_{n \in \mathbb{N}}\left(\frac{1}{n} \sum_{j=1}^{n} a_{j}\right)^{p} \leq\left(p^{\prime}\right)^{p} \sum_{n=1}^{\infty} a_{n}^{p}$.
The constants $\pi / \sin (\pi / p)$ and $\left(p^{\prime}\right)^{p}=(p /(p-1))^{p}$ are the best possible.
At the same time, the continuous versions of the previous inequalities are the
following (see [3, 4]): Let $1<p<\infty$ and $p^{\prime}$ the conjugate of $p$. If $\int_{[0, \infty)}|f|^{p}$ and $\int_{[0, \infty)}|g|^{p^{\prime}}$ are finite, then

$$
\int_{[0, \infty)} \int_{[0, \infty)} \frac{|f(x)||g(y)|}{x+y} d x d y \leq \frac{\pi}{\sin (\pi / p)}\left(\int_{[0, \infty)}|f(x)|^{p} d x\right)^{1 / p}\left(\int_{[0, \infty)}|g(x)|^{p^{\prime}} d x\right)^{1 / p^{\prime}}
$$

and

$$
\int_{[0, \infty)}\left(\frac{1}{x} \int_{[0, x]} f(y) d y\right)^{p} d x \leq\left(\frac{p}{p-1}\right)^{p} \int_{[0, \infty)}|f(x)|^{p} d x
$$

Once again, the constants involved are the best possible.
As usual in harmonic analysis, if $E$ is a measurable subset of $\mathbb{R}^{n}$, then $L^{p}(E)$, $1 \leq p<\infty$, is the Lebesgue space of all measurable functions $f$ such that $\|f\|_{L^{p}(E)}^{p}=$ $\int_{E}|f(x)|^{p} d x$ is finite. Recall that $\left(L^{p}(E),\|\cdot\|_{L^{p}(E)}\right)$ is a Banach space and in the case $E=\mathbb{R}^{n}$, it is denoted $\|\cdot\|_{p}=\|\cdot\|_{L^{p}(E)}$.

Now, consider the operators $H$ and $P$ defined by

$$
H f(x)=\int_{[0, \infty)} \frac{f(t)}{x+t} d t \quad \text { and } \quad P f(x)=\frac{1}{x} \int_{[0, x]} f(t) d t
$$

which naturally arise from the inequalities presented above. Also consider

$$
Q f(x)=\int_{[x, \infty)} \frac{f(t)}{t} d t
$$

being the adjoint operator of $P$ and satisfying

$$
\int_{[0, \infty)}(Q f(x))^{p} d x=\int_{[0, \infty)}\left(\int_{[x, \infty)} \frac{f(t)}{t} d t\right)^{p} d x \leq C \int_{[0, \infty)}(f(x))^{p} d x
$$

for all $f \in L^{p}([0, \infty)), 1<p<\infty$, where $C$ is a positive constant (see [4]). Therefore, $P$ and $Q$ are bounded operators from $L^{p}([0, \infty))$ in itself, that is,

$$
\|P f\|_{L^{p}([0, \infty))} \leq C\|f\|_{L^{p}([0, \infty))} \text { and }\|Q f\|_{L^{p}([0, \infty))} \leq C\|f\|_{L^{p}([0, \infty))} \text { for all } f \in L^{p}([0, \infty)) .
$$

It is immediate that for nonnegative functions $f$,

$$
H f(x) \leq P f(x)+Q f(x) \leq 2 H f(x) \quad \text { forall } x>0 .
$$

Consequently $H$ is a bounded operator on $L^{p}([0, \infty))$, that is,

$$
\|H f\|_{L^{p}([0, \infty))} \leq C\|f\|_{L^{p}([0, \infty))} \quad \text { forall } f \in L^{p}([0, \infty))
$$

It is well known that $P$ is called the Hardy averaging operator and $H$ is called the Hilbert operator. Also, the Calderón operator $S$ is defined by $S=P+Q$, being then a bounded operator from $L^{p}([0, \infty))$ in itself.

We end this section with some of the first and most important direct applications obtained from Hilbert's and Hardy's inequalities.

Theorem 1.1 Let $E$ be the interval $(0,1)$ and $f \in L^{2}(E)$ not null in $E$. Then

$$
\sum_{n=0}^{\infty}\left(\int_{E} x^{n} f(x) d x\right)^{2}<\pi \int_{E} f^{2}(x) d x
$$

and the constant $\pi$ is the best possible. The integrals $\int_{E} x^{n} f(x) d x, n=0,1, \ldots$ are called the moments off in $E$ and are important in many theories.

Theorem 1.2 (Carlema's inequalities) Let $\left\{a_{n}\right\}$ be a sequence of positive numbers and $1<p<\infty$. Then

$$
\sum_{n=1}^{\infty}\left(\frac{1}{n} \sum_{k=1}^{n} a_{k}^{1 / p}\right)^{p}<\left(\frac{p}{p-1}\right)^{p} \sum_{n=1}^{\infty} a_{n} \quad \text { and } \quad \sum_{n=1}^{\infty}\left(\prod_{k=1}^{n} a_{k}\right)^{1 / n}<e \sum_{n=1}^{\infty} a_{n} .
$$

The constants involved are the best possible.
The corresponding integral version for the second inequality of Carlema's inequality is: If $f$ is a positive function belonging to $L^{1}([0, \infty))$, then

$$
\int_{[0, \infty)} \exp \left(\frac{1}{x} \int_{[0, x]} \log f(t) d t\right) d x=\int_{[0, \infty)} e^{P(\log f)(x)} d x<e \int_{[0, \infty)} f(x) d x
$$

where the constant $e$ is the best possible.
Theorem 1.3 Let $1<p \leq 2$ and $p^{\prime}$ the conjugate of $p$. If $L f(s)=\int_{0}^{\infty} f(t) e^{-s t} d t$, i.e. $L f$ is the Laplace transform of $f$, then

$$
\int_{0}^{\infty} L f(s)^{p^{\prime}} d s \leq \frac{2 \pi}{p^{\prime}}\left(\int_{0}^{\infty} f(s)^{p} d s\right)^{p^{\prime} / p} \quad \text { for all } f \in L^{p}([0, \infty))
$$

Therefore $L$ is a bounded operator from $L^{p}([0, \infty))$ into $L^{p^{\prime}}([0, \infty)), 1<p \leq 2$, and $\|L f\|_{p^{\prime}} \leq\left(2 \pi / p^{\prime}\right)^{1 / p^{\prime}}\|f\|_{p}$.

The number of applications and results that arise from Hilbert's and Hardy's inequalities is by now very large and it would be impossible to give a detailed survey of all of them in a reasonable amount of text. We have simply made a very brief introduction about them in this section.

## 2. Calderón weights and $L^{p}$-weighted inequalities

A function $\omega$ defined on $\mathbb{R}^{n}$ is called a weight if it is locally integrable and positive almost everywhere. For a measurable set $E \subset \mathbb{R}^{n},|E|$ denote its Lebesgue measure, $\omega(E)=\int_{E} \omega$, and $E^{c}$ the complement of $E$ in $\mathbb{R}^{n}$. Given a ball $B, t B$ is the ball with the
same center as $B$ and with radius $t$ times as long, and $f_{B}=\frac{1}{|B|} \int_{B} f$. As usual, $\chi_{E}$ denotes the characteristic function of $E$ and $B(x, r)$ denotes a ball centered at $x$ with radius $r$. Also, $C$ denotes a positive constant.

Let $\omega$ be a weight in $\mathbb{R}^{n}$ and $1 \leq p<\infty$. A Lebesgue measurable function $f$ belongs to $L^{p}(\omega)$ if

$$
\|f\|_{L^{p}(\omega)}=\left(\int_{\mathbb{R}^{n}}|f|^{p} \omega\right)^{1 / p}<\infty .
$$

We say that an oprator $T$ is a bounded operator on $L^{p}(\omega)$ if

$$
\|T f\|_{L^{p}(\omega)} \leq C\|f\|_{L^{p}(\omega)}, \quad \text { for all } f \in L^{p}(\omega) .
$$

Given $1<p<\infty$, it is said that $\omega$ is a Calderón weight of class $\mathcal{C}_{p}$, that is $\omega \in \mathcal{C}_{p}$, if the Calderón operator $S$ is bounded on $L^{p}(\omega)$ (see [5]) or, equivalently, if $P$ and $Q$ are both bounded on $L^{p}(\omega)$ (see also [6]). It is well known that the class $\mathcal{C}_{p}$ for $p>1$ is given by the conditions

$$
\begin{aligned}
& M_{p}:\left(\int_{[0, x]} \omega(t) d t\right)^{1 / p}\left(\int_{[x, \infty)} \frac{\omega^{1-p^{\prime}}(t)}{t p^{\prime}} d t\right)^{1 / p^{\prime}} \leq C \quad \text { for all } x>0 \\
& M^{p}:\left(\left(\int_{[x, \infty)} \frac{\omega(t)}{t p} d t\right)^{1 / p}\left(\int_{[0, x]} \omega^{1-p^{\prime}}(t) d t\right)^{1 / p^{\prime}} \leq C \quad \text { for all } x>0 .\right.
\end{aligned}
$$

The Calderón operator plays an important role in the theory of real interpolation and such theory related to Calderón weights is developed in [5]. On the other hand, in [7], the authors considered a maximal operator $N$ on $(0, \infty)$ associated to the basis of open sets of the form $(0, b)$, given by

$$
N f(x)=\sup _{b>x} \frac{1}{b} \int_{[0, b]}|f(t)| d t
$$

for measurable functions $f$. Then, for nonnegative functions $f$, we have

$$
P(x) \leq N f(x) \leq S f(x) \quad \text { forall } x>0
$$

The classes of weights $\omega$ associated to the boundedness of $N$ on $L^{p}(\omega)$ are those that satisfy the Muckenhoupt- $A_{p}$ condition, $1 \leq p<\infty$, only for the sets of the form $(0, b)$. These classes are denoted by $A_{p, 0}$ and defined as follows:

$$
\begin{array}{ll}
A_{1,0}: & N \omega(x) \leq C \omega(x) \text { for almost all } x>0 ; \\
A_{p, 0}: & \left(\frac{1}{x} \int_{[0, x]} \omega\right)\left(\frac{1}{x} \int_{[0, x]} \omega^{1-p^{\prime}}\right)^{p-1} \leq C \quad \text { for all } x>0, \text { where } 1<p<\infty .
\end{array}
$$

Then, in [7] is proved that $N$ and $S$ are bounded operators on $L^{p}(\omega)$ if and only if $\omega \in A_{p, 0}$ for $1<p<\infty$. This result implies, in particular, that the classes of weights $C_{p}$ and $A_{p, 0}$ coincide for $1<p<\infty$.

Taking into account these results it is natural to wonder for the action of the Calderón and Hilbert operators over suitable spaces such as BMO or Lipschitz spaces. Also, another interesting question is: which are, in these cases, the Calderón weights in order to obtain weighted inequalities between these spaces?

These problems were treated for instance in the case of the fractional integral operator in [8, 9], which have been the main motivation for the article [10] and for the development of the following sections.

## 3. The $\boldsymbol{n}$-dimensional Calderón and Hilbert operators

For $0 \leq \alpha<n, f$ a Lebesgue measurable function and $x \in \mathbb{R}^{n}, x \neq 0$, the general $n$-dimensional Calderón and Hilbert operators are defined by

$$
S_{\alpha} f(x)=P_{\alpha} f(x)+Q_{a} f(x) \quad \text { and } \quad H_{\alpha} f(x)=\int_{\mathbb{R}^{n}} \frac{f(y)}{(|x|+|y|)^{n-\alpha}} d y
$$

where $P_{\alpha} f(x)=\frac{1}{|x|^{n-\alpha}} \int_{|y| \leq|x|} f(y) d y$ and $Q_{\alpha} f(x)=\int_{|y|>|x|} \frac{f(y)}{|y|^{n-\alpha}} d y$.
Again, it is immediate that for nonnegative functions $f$, the following pointwise inequalities hold

$$
\begin{equation*}
H_{\alpha} f(x) \leq S_{\alpha} f(x) \leq 2^{n-\alpha} H_{\alpha} f(x) \tag{1}
\end{equation*}
$$

and consequently, all weighted $-L^{p}$ inequalities obtained for $S$ are true for $H$ and reciprocally.

In spite of the punctual comparison (1), we will show in Section 4 that the results obtained for $S_{\alpha}$ and $H_{\alpha}$ are not analogous when the $B M O^{\gamma}$ and Lipschitz spaces are involved.

Both operators $S_{\alpha}$ and $H_{\alpha}$ appear in several different contexts and applications, see for instance [4, 11-17].

Next, we introduce the spaces of functions and the classes of weights which appear in our main results.

Recall that a measurable function $f$ defined on $E \subset \mathbb{R}^{n}$ is said to be essentially bounded provided there is some $M \geq 0$, called an essential upper bound for $f$, for which $|f(x)| \leq M$ for almost all $x \in E$. As usual, the class of all functions that are essentially bounded on $E$ is denoted by $L^{\infty}(E)$ and $\|f\|_{\infty}$ is the infimum of the essential upper bounds for $f \in L^{\infty}(E)$. Then, $\left(L^{\infty}(E),\|\cdot\|_{\infty}\right)$ is a Banach space.

Now, a Lebesgue measurable function $f$ belongs to $L^{\infty}(\omega)$ if $\|f \omega\|_{\infty}<\infty$.
Also recall that $L_{l o c}^{1}\left(\mathbb{R}^{n}\right)$ denotes the space of locally integrable functions $f$ satisfying that $\left\|f \chi_{B}\right\|_{1}$ is finite for every ball $B \subset \mathbb{R}^{n}$.

Definition 3.2. Let $\omega$ be a weight in $\mathbb{R}^{n}$ and $0 \leq \gamma<1 / n$. A locally integrable function $f$ belongs to $B M O^{\gamma}(\omega)$ if there exists a constant $C$ such that for every ball $B \subset \mathbb{R}^{n}$,

$$
\begin{equation*}
\frac{1}{\omega(B)|B|^{\gamma}} \int_{B}\left|f-f_{B}\right| \leq C . \tag{2}
\end{equation*}
$$

The seminorm of $f \in B M O^{\gamma}(\omega),\|f\|_{B M O^{\gamma}(\omega)}$, is the infimum of all such $C$.

Definition 3.4. Let $\omega$ be a weight in $\mathbb{R}^{n}$ and $0 \leq \gamma<1 / n$. A locally integrable function $f$ belongs to $B M_{0}^{\gamma}(\omega)$ if there exists a constant $C$ such that

$$
\begin{equation*}
\frac{1}{\omega(B)|B|^{\gamma}} \int_{B}|f| \leq C \tag{3}
\end{equation*}
$$

for every ball $B \subset \mathbb{R}^{n}$ centered at the origin.
The norm of $f \in B M_{0}^{\gamma}(\omega)$, denoted by $\|f\|_{B M_{0}^{\gamma}(\omega)}$, is the infimum of all such $C$. We will denote by $B M_{0}(\omega)=B M_{0}^{0}(\omega)$.

Observe that with these definitions the space $B M O^{0}(\omega)$ is the weighted version of $B M O$ introduced by Muckenhoupt and Wheeden in [18]. Also, the family of spaces $B M O^{\gamma}(\omega)$ is contained in the family of weighted Lipschitz spaces $\mathcal{I}_{\omega}(\gamma)$ defined and studied in [8], and $B M O^{\gamma}(\omega)$ for $\omega \equiv 1$ is the well known Lipschitz integral space. Furthermore, we note that $L^{\infty}\left(\omega^{-1}\right) \subset B M_{0}(\omega) \cap B M O(\omega)$.

Given $p>1$, it is known that a weight $\omega$ satisfies the reverse Hölder inequality with exponent $p$, denoted by $\omega \in R H(p)$, if

$$
\begin{equation*}
\left(\frac{1}{|B|} \int_{B} \omega^{p}\right)^{1 / p} \leq C \frac{1}{|B|} \int_{B} \omega \tag{4}
\end{equation*}
$$

for all balls $B \subset \mathbb{R}^{n}$ and some constant $C$.
Definition 3.7. Given $p>1$, a weight $\omega$ belongs to $R H_{0}(p)$ if it satisfies (4) but only for balls centered at the origin.

Definition 3.8. A weight $\omega$ belongs to $D_{0}$ if it satisfies the doubling condition $\omega(2 B) \leq C \omega(B)$ for every ball $B \subset \mathbb{R}^{n}$ centered at the origin and some constant $C$.

Definition 3.9. Let $\eta \geq 1$, a weight $\omega$ belongs to $D_{\eta}$ if it satisfies the doubling condition

$$
\frac{\omega(2 B(x,|x|+r))}{|B(x,|x|+r)|^{\eta}} \leq C \frac{\omega(B(x, r))}{|B(x, r)|^{\eta}}
$$

every ball $B(x, r) \subset \mathbb{R}^{n}$ and some constant $C$.
It is immediate that $D_{\eta} \subset D_{0}$ for all $\eta$, and $D_{\eta}$ is increasing in $\eta$. It is well known that each weight in the Muckenhoupt class $A_{\infty}$ is in $R H(p) \cap D_{\eta}$ for some $p$ and for some $\eta$, see for instance [19]. On the other hand, there exist weights belonging to $D_{\eta}$ for some $\eta$, such that it is not in $A_{\infty}$, see [20].

Also, we observe the following property that we will use along this chapter. If $\omega \in D_{\eta}$ there exists a constant $C$ such that

$$
\begin{equation*}
\omega(B) \leq C \omega\left(B \backslash \frac{1}{2} B\right) \tag{5}
\end{equation*}
$$

for every ball $B \subset \mathbb{R}^{n}$ centered at the origin.
Definition 3.11. Let $0 \leq \alpha<n$ and $1<p<\infty$. A weight $\omega$ belongs to $\mathrm{H}_{0}(\alpha, p)$ if there exists a constant $C$ such that

$$
\begin{equation*}
\left(\int_{B^{c}} \frac{\omega^{p^{\prime}}(y)}{|y|^{(n-\alpha+1)} p^{\prime}} d y\right)^{1 / p^{\prime}} \leq C \frac{\omega(B)}{|B|^{1+1 / p-\alpha / n+1 / n}} \tag{6}
\end{equation*}
$$

for every ball $B \subset \mathbb{R}^{n}$ centered at the origin.
A weight $\omega$ belongs to $\mathrm{H}_{0}(\alpha, \infty)$ if there exists a constant $C$ such that

$$
\begin{equation*}
\int_{B^{c}} \frac{\omega(y)}{|y|^{n-\alpha+1}} d y \leq C \frac{\omega(B)}{|B|^{1-\alpha / n+1 / n}} \tag{7}
\end{equation*}
$$

for every ball $B \subset \mathbb{R}^{n}$ centered at the origin.
The classes of weights $\mathrm{H}_{0}(\alpha, p)$ and $\mathrm{H}_{0}(\alpha, \infty)$ satisfying (6) and (7) respectively but for all ball $B \subset \mathbb{R}^{n}$, were introduced and studied in [8].

## 4. Weighted Lebesgue and $B M O^{\gamma}$ norm inequalities for $S_{\alpha}$ and $H_{\alpha}$

Before beginning our study of the generalized Calderón operator, we notice that $S_{\alpha} f$ can be identically infinite for some functions $f$ belonging to $L^{p}\left(\omega^{-p}\right)$ or $B M_{0}^{\gamma}(\omega)$. For example, for $\omega \equiv 1$ and $\alpha>0$, if $f(x)=|x|^{-\alpha} \chi_{B^{c}(0,1)}(x)$ and $n / \alpha<p$, then $f \in L^{p}\left(\omega^{-p}\right)$ but $S_{\alpha} f \equiv \infty$. For the case $n / \alpha=p$, if $g(x)=|x|^{-\alpha}(\log |x|)^{-(1+1 / p) / 2} \chi_{B^{c}(0,2)}(x)$, then $g \in L^{p}\left(\omega^{-p}\right)$ but $S_{\alpha} g \equiv \infty$. Also, if $h(x)=\chi_{B^{c}(0,1)}(x)$, then $h \in B M_{0}^{\gamma}(\omega)$ but $S_{\alpha} h \equiv \infty$ for all $0 \leq \alpha<n$. However, in Lemma 4.7 we will show that if $f$ belongs to $L^{p}\left(\omega^{-p}\right) \cup B M_{0}^{\gamma}(\omega)$ and $S_{\alpha} f(x)$ is finite for some $x \neq 0$, then $S_{\alpha} f$ is finite on $\mathbb{R}^{n} \backslash\{0\}$. This also happens for the generalized Hilbert operator since the comparison (1).

Therefore, throughout the following sections we shall consider $S_{\alpha}$ and $H_{\alpha}$ defined on functions $f$ belonging to $L^{p}\left(\omega^{-p}\right)$ or $B M_{0}^{\gamma}(\omega)$ such that $S_{\alpha} f$ and $H_{a} f$ are finite for some $x \neq 0$.

Also, note that $S_{\alpha} f$ is finite on $\mathbb{R}^{n} \backslash\{0\}$ for all compactly supported functions $f \in L^{\infty}\left(\omega^{-1}\right)$, and the same holds for $H_{\alpha} f$. These functions belongs to $L^{p}\left(\omega^{-p}\right)$ and those such that zero is not in their support belongs to $B M_{0}^{\gamma}(\omega)$.

The operator $P$ is naturally bounded from $B M_{0}$ into $L^{\infty}$ and analogously, $Q$ is naturally bounded from $B M_{0}$ into $B M O$ (see Proposition 3.5 in [13]). So, immediately the Calderón operator is bounded from $B M_{0}$ into $B M O$. This natural boundedness is our motivation in order to consider the $B M_{0}^{\gamma}(\omega)$ spaces and obtain Theorems 1.5 and 1.7. Likewise, since $L^{\infty}\left(\omega^{-1}\right) \subset B M_{0}(\omega)$, we get Corollaries 4.1 and 4.2.

We now state the main results of this chapter.
Theorem 1.4 Suppose $\alpha>0, n / \alpha \leq p<n /(\alpha-1)^{+}, \eta=1+1 / n+1 / p-\alpha / n$ and $\delta=\alpha / n-1 / p$. The operator $S_{\alpha}$ is bounded from $L^{p}\left(\omega^{-p}\right)$ into $B M O^{\delta}(\omega)$ and $\omega^{p^{\prime}} \in D_{0}$ if and only if $\omega \in R H_{0}\left(p^{\prime}\right) \cap D_{\eta}$.

Theorem 1.5 Suppose $0 \leq \alpha<1,0 \leq \gamma<1 / n-\alpha / n, \eta=1+1 / n-\alpha / n-\gamma$ and $\delta=$ $\alpha / n+\gamma$. The operator $S_{\alpha}$ is bounded from $B M_{0}^{\gamma}(\omega)$ into $B M O^{\delta}(\omega)$ and $\omega \in D_{0}$ if and only if $\omega \in D_{\eta}$.

Corollary 4.1. Let $\eta=1+1 / n$. Then $S$ is bounded from $L^{\infty}\left(\omega^{-1}\right)$ into $B M O(\omega)$ and $\omega \in D_{0}$ if and only if $\omega \in D_{\eta}$.

Theorem 1.6 Suppose $\alpha>0, n / \alpha \leq p<n /(\alpha-1)^{+}, \eta=1+1 / n+1 / p-\alpha / n$ and $\delta=\alpha / n-1 / p$. The operator $H_{\alpha}$ is bounded from $L^{p}\left(\omega^{-p}\right)$ into $B M O^{\delta}(\omega)$ if and only if $\omega \in \mathrm{H}_{0}(\alpha, p) \cap R H_{0}\left(p^{\prime}\right) \cap D_{\eta}$.

Theorem 1.7 Suppose $0 \leq \alpha<1,0 \leq \gamma<1 / n-\alpha / n, \eta=1+1 / n-\alpha / n-\gamma$ and $\delta=$ $\alpha / n+\gamma$. The operator $H_{\alpha}$ is bounded from $B M_{0}^{\gamma}(\omega)$ into $B M O^{\delta}(\omega)$ if and only if $\omega \in \mathrm{H}_{0}(\alpha+n \gamma, \infty) \cap D_{\eta}$.

Corollary 4.2. Let $\eta=1+1 / n$. Then $H$ is bounded from $L^{\infty}\left(\omega^{-1}\right)$ into $B M O(\omega)$ if and only if $\omega \in \mathrm{H}_{0}(0, \infty) \cap D_{\eta}$.

Remark 4.3. It is classic the study of the boundedness of operators between $L^{\infty}$ and $B M O$ spaces. In [10], the results obtained in Corollaries 4.1 and 4.2 are originals, even in the unweighted case for $H$. The unweighted case for $S$ is contained in Proposition 3.5 of [13].

Remark 4.4. The limit case $p=\infty\left(p^{\prime}=1\right)$ of Theorem 1.4 is contained in Theorem 1.5 with $\gamma=0$, since the hypotheses on the weights coincide. This also is true to Theorems 1.6 and 1.7.

Let $\alpha, p$ and $\eta$ be as in Theorems 1.4 and 1.6. It is not difficult to show that if $\omega^{p^{\prime}} \in A_{1,0}$ then $\omega \in \mathrm{H}_{0}(\alpha, p) \cap R H_{0}\left(p^{\prime}\right) \cap D_{\eta}$. Also, if $\omega(x)=|x|^{\beta}$ with $\beta \in(0,1+n / p-\alpha)$, then $\omega^{p^{\prime}} \notin A_{1,0}$ but $\omega \in \mathrm{H}_{0}(\alpha, p) \cap R H_{0}\left(p^{\prime}\right) \cap D_{\eta}$. Furthermore, if $\omega(x)=|x|^{\beta}$ with $\beta=1+n / p-\alpha$, then $\omega \in R H_{0}\left(p^{\prime}\right) \cap D_{\eta}$ but $\omega \notin \mathrm{H}_{0}(\alpha, p)$. Now, if in addition $0<\alpha<1$ and $p^{\prime}>n /(1-\alpha)$, we have that if $\omega^{p^{\prime}} \in A_{p^{\prime}+1,0}$ then $\omega \in \mathrm{H}_{0}(\alpha, p) \cap R H_{0}\left(p^{\prime}\right) \cap D_{\eta}$. In fact, the $\mathrm{H}_{0}(\alpha, p)$-condition is obtained directly from the $A_{p^{\prime}+1,0}$-condition, and by Hölder inequality we have that

$$
\begin{aligned}
\left(\frac{\omega^{p^{\prime}}\left(B\left(0,2\left(\left|x_{0}\right|+r\right)\right)\right)}{\left|B\left(0,2\left(\left|x_{0}\right|+r\right)\right)\right|}\right)^{1 / p^{\prime}} & \leq C \frac{\left|B\left(0,2\left(\left|x_{0}\right|+r\right)\right)\right|}{\omega^{-1}\left(B\left(x_{0}, r\right)\right)} \leq C\left(\frac{\left|x_{0}\right|+r}{r}\right)^{n} \frac{\omega\left(B\left(x_{0}, r\right)\right)}{\left|B\left(x_{0}, r\right)\right|} \\
& \leq C\left(\frac{\left|x_{0}\right|+r}{r}\right)^{1-\alpha+n / p} \frac{\omega\left(B\left(x_{0}, r\right)\right)}{\left|B\left(x_{0}, r\right)\right|}
\end{aligned}
$$

for all balls $B\left(x_{0}, r\right) \subset \mathbb{R}^{n}$. Thus, the $R H_{0}\left(p^{\prime}\right)$ and $D_{\eta}$ conditions follow from the last expression.

On the other hand, suppose that $\alpha, \gamma$ and $\eta$ be as in Theorems 1.5 and 1.7. If $\omega \in A_{1,0}$ then $\omega \in \mathrm{H}_{0}(\alpha+n \gamma, \infty) \cap D_{\eta}$. Also, if $\omega(x)=|x|^{\beta}$ with $\beta \in(0,1-\alpha-n \gamma)$, then $\omega \notin A_{1,0}$ but $\omega \in \mathrm{H}_{0}(\alpha+n \gamma, \infty) \cap D_{\eta}$. Finally, if $\omega(x)=|x|^{\beta}$ with $\beta=1-\alpha-n \gamma$, then $\omega \in D_{\eta}$ but $\omega \notin \mathrm{H}_{0}(\alpha+n \gamma, \infty)$.

We shall denote by $A(x, r, R)$ with $0<r<R$ the annulus centered at $x$ with radii $r$ and $R$, and by $C$ and $c$ positive constants not necessarily the same at each occurrence.

Before proceeding to the proofs of the main theorems we give some previous lemmas.

Suppose that $1<p<\infty$ and $\omega \in R H_{0}\left(p^{\prime}\right)$, then it is easy to see that there exists $C$ such that

$$
\begin{equation*}
\int_{B}|f| \leq C \frac{\omega(B)}{|B|^{1} / p}\|f\|_{L^{p}}\left(\omega^{-p}\right) \tag{8}
\end{equation*}
$$

for all $f \in L^{p}\left(\omega^{-p}\right)$ and for every ball $B \subset \mathbb{R}^{n}$ centered at the origin.
Lemma 4.6. (i) Let $0<\alpha<n$ and $1<p<\infty$. If $\omega \in \mathrm{H}_{0}(\alpha, p)$ then there exists $C$ such that

$$
\int_{B^{c}} \frac{|f(y)|}{|y|^{n-\alpha+1}} d y \leq C \frac{\omega(B)}{|B|^{1+1 / p-\alpha / n+1 / n}}\|f\|_{L^{p}\left(\omega^{-p}\right)}
$$

for all $f \in L^{p}\left(\omega^{-p}\right)$ and for every ball $B \subset \mathbb{R}^{n}$ centered at the origin.
(ii) Let $0 \leq \alpha<1,0 \leq \gamma<1 / n-\alpha / n$ and $\eta=1+1 / n-\alpha / n-\gamma$. If $\omega \in \mathrm{H}_{0}(\alpha+n \gamma, \infty) \cap D_{\eta}$ then there exists $C$ such that

$$
\int_{B^{c}} \frac{|f(y)|}{|y|^{n-\alpha+1}} d y \leq C \frac{\omega(B)}{|B|^{\eta}}\|f\|_{B M_{0}^{\gamma}(\omega)}
$$

for all $f \in B M_{0}^{\gamma}(\omega)$ and for every ball $B \subset \mathbb{R}^{n}$ centered at the origin.
Proof: The part $(i)$ is immediate from Hölder's inequality and Definition 3.11. For (ii), since the hypothesis on $\omega$ and (3.10), for $B=B(0, r)$ we have

$$
\begin{aligned}
\int_{B^{c}} \frac{|f(y)|}{|y|^{n-\alpha+1}} d y & \leq C \sum_{k=0}^{\infty} \frac{1}{\left(2^{k} r\right)^{n-\alpha+1}} \int_{2^{k} r \leq y \mid<2^{k+1} r}|f(y)| d y \\
& \leq C\|f\|_{B M_{0}^{\gamma}(\omega)} \sum_{k=0}^{\infty} \frac{\omega\left(B\left(0,2^{k+1} r\right)\right)}{\left(2^{k} r\right)^{n-\alpha+1-n \gamma}} \\
& \leq C\|f\|_{B M_{0}^{y}(\omega)} \sum_{k=0}^{\infty} \frac{\omega\left(B\left(0,2^{k+1} r\right) \backslash B\left(0,2^{k} r\right)\right)}{\left(2^{k} r\right)^{n-(\alpha+n \gamma)+1}} \\
& \leq C\|f\|_{B M_{0}^{\gamma}(\omega)} \frac{\omega(B)}{|B|^{\eta}} .
\end{aligned}
$$

Lemma 4.7. (i) Let $\alpha>0,1<p<\infty$ and $\omega \in R H_{0}\left(p^{\prime}\right)$. If $f \in L^{p}\left(\omega^{-p}\right)$ and there exists $x \neq 0$ such that $S_{\alpha} f(x)$ is finite, then $S_{a f}$ is finite on $\mathbb{R}^{n} \backslash\{0\}$ and $S_{a} f \in L_{l o c}^{1}\left(\mathbb{R}^{n}\right)$. The claim also holds for $H_{\alpha}$.
(ii) Let $\omega \in D_{\eta}$. Iff $\in B M_{0}^{\gamma}(\omega)$ and there exists $x \neq 0$ such that $S_{\alpha} f(x)$ is finite, then $S_{\alpha} f$ is finite on $\mathbb{R}^{n} \backslash\{0\}$ and $S_{\alpha} f \in L_{l o c}^{1}\left(\mathbb{R}^{n}\right)$. The claim also holds for $H_{\alpha}$.

Proof: Since (3.1) we will only consider the operator $S_{\alpha}$. Suppose $f$ is a nonnegative function in $L_{l o c}^{1}\left(\mathbb{R}^{n}\right)$ such that $S_{\alpha} f\left(x_{0}\right)<\infty$ for some $x_{0} \neq 0$. Then $Q_{\alpha} f(x)<\infty$ for $|x| \geq\left|x_{0}\right|$, and if $0<|x|<\left|x_{0}\right|$ then

$$
Q_{a} f(x) \leq \frac{1}{|x|^{n-\alpha}} \int_{|x|<|y|<\left|x_{0}\right|} f(y) d y+Q_{a} f\left(x_{0}\right)<\infty .
$$

Furthermore, since

$$
\int_{B(0, r)}\left(Q_{\alpha} f(x)-Q_{\alpha} f(\nu)\right) d x \leq \int_{B(0, r)} f(y) r^{\alpha} d y<\infty,
$$

where $|\nu|=r$, then $Q_{\alpha} f \in L_{l o c}^{1}\left(\mathbb{R}^{n}\right)$.
If $\alpha>0$ it is immediate that $P_{\alpha} f \in L_{l o c}^{1}\left(\mathbb{R}^{n}\right)$. Therefore, (i) follows from (4.5). For (ii) it remains to show that $P_{\alpha} f \in L_{l o c}^{1}\left(\mathbb{R}^{n}\right)$ in the case $\alpha=0$. Let $B_{j}=B\left(0,2^{-j} r\right), j=$ $0,1, \ldots$, by (3.10) we have

$$
\begin{aligned}
\int_{B_{0}} \frac{1}{|x|^{n}} \int_{B(0,|x|)} f(y) d y d x & \leq C\|f\|_{B M_{0}^{\gamma}(\omega)} \int_{B_{0}} \frac{\omega(B(0,|x|))}{|x|^{n-n \gamma}} d x \\
& \leq C\|f\|_{B M_{0}^{\gamma}(\omega)} \sum_{j=0}^{\infty} \frac{r^{n \gamma-n}}{2^{j(n \gamma-n)}} \int_{B_{j} \backslash B_{j+1}} \omega\left(B_{j}\right) d x \\
& \leq C\|f\|_{B M_{0}^{\gamma}(\omega)} r^{r^{n \gamma}} \sum_{j=0}^{\infty} \frac{\omega\left(B_{j} \backslash B_{j+1}\right)}{2^{j n \gamma}} \\
& \leq C\|f\|_{B M_{0}^{\gamma}(\omega)} r^{r^{n \gamma}} \omega\left(B_{0}\right) .
\end{aligned}
$$

Proof of Theorem 1.4: We begin showing the sufficient condition. Let $B=$ $B\left(x_{0}, r\right)$. If $x_{0}=0$, let $u=r e_{1} / 2$ and $v=3 r e_{1} / 4$, where $e_{1}=(1, \ldots, 0)$. If $x_{0} \neq 0$, let $u=\left(\left|x_{0}\right|+r / 2\right) x_{0} /\left|x_{0}\right|$ and $v=\left(\left|x_{0}\right|+3 r / 4\right) x_{0} /\left|x_{0}\right|$. Thus, we consider the following two regions

$$
\begin{align*}
& U=B(u, r / 8) \cap\left\{u+h: \operatorname{sign}\left(u_{i}\right)=\operatorname{sign}\left(h_{i}\right) i=1, \ldots, n\right\}, \\
& V=B(v, r / 4) \cap\left\{v+h: \operatorname{sign}\left(v_{i}\right)=\operatorname{sign}\left(h_{i}\right) i=1, \ldots, n\right\}, \tag{9}
\end{align*}
$$

where $u=\left(u_{1}, \ldots, u_{n}\right), v=\left(v_{1}, \ldots, v_{n}\right)$ and $h=\left(h_{1}, \ldots, h_{n}\right)$. In the case $u_{i}=0$ for some $i$, we choose $h_{i}>0$. Clearly, we have the estimates $\operatorname{dist}(U, V)=C r$,

$$
|U|=\frac{1}{2^{n}}|B(u, r / 8)|=C|B| \quad \text { and } \quad|V|=\frac{1}{2^{n}}|B(v, r / 4)|=C|B| .
$$

Let $f$ a nonnegative function in $L^{p}\left(\omega^{-p}\right)$ such that $\operatorname{supp}(f) \subset B\left(0,\left|x_{0}\right|+r / 2\right)$, where $\operatorname{supp}(f)$ is the closure of the set $\{x: f(x) \neq 0\}$. Then

$$
\begin{aligned}
\left\|S_{\alpha} f\right\|_{B M O^{\delta}(\omega)} & \geq \frac{C}{\omega(B)|B|^{1+\delta}} \int_{B} \int_{B}\left|S_{\alpha} f(x)-S_{\alpha} f(z)\right| d z d x \\
& \geq \frac{C}{\omega(B)|B|^{1+\delta}} \int_{U} \int_{V}\left|\left(\frac{1}{|x|^{n-\alpha}}-\frac{1}{|z|^{n-\alpha}}\right) \int_{B\left(0,\left|x_{0}\right|+r / 2\right)} f(y) d y\right| d z d x .
\end{aligned}
$$

Note that, for $x \in U$ and $z \in V$ we have $\frac{1}{|x|^{n-\alpha}}-\frac{1}{|z|^{n-\alpha}} \geq C \frac{r}{\left(\left|x_{0}\right|+r\right)^{n-\alpha+1}}$. Then

$$
\begin{equation*}
\left\|S_{\alpha} f\right\|_{B M O^{\delta}(\omega)} \geq \frac{C r^{n+1}}{\omega(B)|B|^{\delta}\left(\left|x_{0}\right|+r\right)^{n-\alpha+1}} \int_{B\left(0,\left|x_{0}\right|+r / 2\right)} f(y) d y \tag{10}
\end{equation*}
$$

Thus, taking $f(y)=\omega^{p^{\prime}}(y) \chi_{B\left(0,\left|x_{0}\right|+r / 2\right)}(y)$ in (10) and since the boundedness of $S_{\alpha}$ and $\omega^{p^{\prime}} \in D_{0}$, we have

$$
\left(\frac{\omega^{p^{\prime}}\left(B\left(0,\left|x_{0}\right|+r\right)\right.}{\left|B\left(0,\left|x_{0}\right|+r\right)\right|}\right)^{1 / p^{\prime}} \leq C\left(\frac{\left|x_{0}\right|+r}{r}\right)^{1-\alpha+n / p} \frac{\omega(B)}{|B|}
$$

Taking $x_{0}=0$ in the last expression, we have that $\omega \in R H_{0}\left(p^{\prime}\right)$. Then, applying the Hölder's inequality, we obtain that $\omega$ satisfies the desired condition $D_{\eta}$.

Now, let us show the necessary condition. Let $f \in L^{p}\left(\omega^{-p}\right)$ such that $S_{\alpha} f(x)$ is finite for some $x \neq 0$ and let $\omega \in R H_{0}\left(p^{\prime}\right) \cap D_{\eta}$. It is immediate that $\omega^{p^{\prime}} \in D_{0}$. Thus $S_{\alpha} f \in L_{l o c}^{1}\left(\mathbb{R}^{n}\right)$ by $(i)$ of Lemma 4.7. First, we consider $B=B(0, r), x \in B$ and $x \neq 0$. Let $\nu$ be such that $|\nu|=r$, and let

$$
K_{\nu}(x, y)=\min \left\{1, \frac{|y|^{n-\alpha}}{|x|^{n-\alpha}}\right\}-\min \left\{1, \frac{|y|^{n-\alpha}}{|\nu|^{n-\alpha}}\right\} .
$$

Then, since $K_{\nu}(x, y)=0$ for $|y|>|\nu|$, we have

$$
\begin{equation*}
S_{\alpha} f(x)-S_{\alpha} f(\nu)=\int_{|y| \leq|\nu|} K_{\nu}(x, y) \frac{f(y)}{|y|^{n-\alpha}} d y . \tag{11}
\end{equation*}
$$

If $|y| \leq|\nu|$ then $K_{\nu}(x, y) \geq 0$, so

$$
\begin{align*}
\frac{1}{\omega(B)} \int_{B}\left|S_{a} f(x)-S_{\alpha} f(\nu)\right| d x & \leq \frac{1}{\omega(B)} \int_{B} \int_{B} K_{\nu}(x, y) \frac{|f(y)|}{|y|^{n-\alpha}} d y d x \\
& =\frac{1}{\omega(B)} \int_{B} \int_{\varphi|\leq|x|} K_{\nu}(x, y) \frac{|f(y)|}{|y|^{n-\alpha}} d y d x+\frac{1}{\omega(B)} \int_{B} \int_{|x|<\langle y| \leq r} K_{\nu}(x, y) \frac{|f(y)|}{|y|^{n-\alpha}} d y d x . \tag{12}
\end{align*}
$$

Now we estimate each term in (12).
If $|y| \leq|x|$ then $K_{\nu}(x, y) \leq|y|^{n-\alpha}|x|^{-(n-\alpha)}$. So, by (8) we have

$$
\begin{aligned}
\frac{1}{\omega(B)} \int_{B} \int_{|y| \leq|x|} K_{\nu}(x, y) \frac{|f(y)|}{|y|^{n-\alpha}} d y d x & \leq \frac{1}{\omega(B)} \int_{B} \frac{1}{|x|^{n-\alpha}} \int_{B}|f(y)| d y d x \\
& \leq\left. C\left|f \|_{L^{p}\left(\omega^{-p}\right)}\right| B\right|^{\delta} .
\end{aligned}
$$

For the second term, since $0 \leq K_{\nu}(x, y) \leq 1$ and (8), we have

$$
\begin{align*}
\frac{1}{\omega(B)} \int_{B} \int_{|x|<|y| \leq r} K_{\nu}(x, y) \frac{|f(y)|}{|y|^{n-\alpha}} d y d x & \leq \frac{1}{\omega(B)} \int_{B} \frac{1}{|x|^{n-\alpha}} \int_{|x|<|y| \leq r}|f(y)| d y d x \\
& \leq \frac{C}{\omega(B)} \int_{B}|f(y) \| y|^{\alpha} d y  \tag{13}\\
& \leq C\|f\|_{L^{p}\left(\omega^{-p}\right)}|B|^{\delta} .
\end{align*}
$$

Then, by (12) and (13), we have proved

$$
\begin{equation*}
\frac{1}{\omega(B)|B|^{\delta}} \int_{B}\left|S_{\alpha} f(x)-S_{\alpha} f(\nu)\right| d x \leq C\|f\|_{L^{p}\left(\omega^{-p}\right)}, \tag{14}
\end{equation*}
$$

for every ball $B$ centered at the origin.
We now consider $B=B\left(x_{0}, r\right)$ with $r<\left|x_{0}\right| / 8$. By (14) it is enough to consider only these balls $B$. Let $x \in B$ and $\nu=\left(\left|x_{0}\right|+r\right) x_{0} /\left|x_{0}\right|$. In the same way as (11), we have

$$
S_{\alpha} f(x)-S_{\alpha} f(\nu)=\int_{|v| \leq|\nu|} K_{\nu}(x, y) \frac{f(y)}{|y|^{n-\alpha}} d y .
$$

Now, we note that if $|y| \leq|\nu|$ then $K_{\nu}(x, y) \geq 0$. Applying the mean value theorem and using $|\nu| \sim|x|$, then

$$
\begin{equation*}
K_{\nu}(x, y) \leq \frac{|y|^{n-\alpha}}{|x|^{n-\alpha}}-\frac{|y|^{n-\alpha}}{|\nu|^{n-\alpha}} \leq C \frac{r|y|^{n-\alpha}}{|\nu|^{n-\alpha+1}} . \tag{15}
\end{equation*}
$$

Thus, by (8) and $\omega \in D_{\eta}$, we have

$$
\begin{align*}
\frac{1}{\omega(B)} \int_{B}\left|S_{\alpha} f(x)-S_{\alpha} f(\nu)\right| d x & \left.\leq C \frac{r}{\omega(B)|\nu|^{n-\alpha+1}} \int_{B} \int_{|v| \leq|\nu|} f(y) \right\rvert\, d y d x \\
& \leq C\|f\|_{L^{p}\left(\omega^{-p}\right)} \frac{r^{n+1}|\nu|^{n-\alpha+1+n / p}}{} \frac{\omega(B(0,|\nu|))}{\omega(B)}  \tag{16}\\
& \leq C\|f\|_{L^{p}\left(\omega^{-p}\right)}|B|^{\delta} .
\end{align*}
$$

Therefore, (14) and (16) complete the proof of the theorem.
Proof of Theorem 1.5: We begin showing the sufficient condition. Let $B=B\left(x_{0}, r\right)$ and let $u, v, U$ and $V$ as in (9) of the proof of Theorem 1.4. Then, we again have

$$
\begin{equation*}
\left\|S_{\alpha} f\right\|_{B M O^{\delta}(\omega)} \geq \frac{C r^{n+1}}{\omega(B)|B|^{\delta}\left(\left|x_{0}\right|+r\right)^{n-\alpha+1}} \int_{B\left(0,\left|x_{0}\right|+r / 2\right)} f(y) d y \tag{17}
\end{equation*}
$$

for every nonnegative function $f$ in $B M_{0}^{\gamma}(\omega)$ such that $\operatorname{supp}(f) \subset B\left(0,\left|x_{0}\right|+r / 2\right)$. Now, if $\gamma=0$ we take $f(y)=\omega(y) \chi_{B\left(0,\left|x_{0}\right|+r / 2\right)}(y)$ in (17) and since $\|f\|_{B M_{0}^{\%}(\omega)} \leq 1$, the boundedness of $S_{\alpha}$ and $\omega \in D_{0}$, we have $\omega \in D_{\eta}$.

If $\gamma>0$, let $f(y)=P_{n \gamma}\left(\omega \chi_{B\left(0,\left|x_{0}\right|+r / 2\right)}\right)(y)$, then $\|f\|_{B M_{0}^{\gamma}(\omega)} \leq C$ and

$$
\begin{align*}
\int_{B\left(0,\left|x_{0}\right|+r / 2\right)} f(y) d y & =C \int_{B\left(0,\left|x_{0}\right|+r / 2\right)} \omega(t)\left(\left(\left|x_{0}\right|+r / 2\right)^{n \gamma}-|t|^{n \gamma}\right) d t  \tag{18}\\
& \geq C\left(\left|x_{0}\right|+r\right)^{n \gamma} \omega\left(B\left(0,\left(\left|x_{0}\right|+r / 2\right) / 2\right)\right) .
\end{align*}
$$

Therefore, using this function $f$ in (17), the boundedness of $S_{\alpha}$, (18) and $\omega \in D_{0}$, we have $\omega \in D_{\eta}$.

Now, let us show the necessary condition. Let $f \in B M_{0}^{\gamma}(\omega)$ such that $S_{\alpha} f(x)$ is finite for some $x \neq 0$ and let $\omega \in D_{\eta}$. Thus $S_{\alpha} f \in L_{l o c}^{1}\left(\mathbb{R}^{n}\right)$ by (ii) of Lemma 4.7. We begin considering $B=B(0, r), x \in B$ and $x \neq 0$. Let $\nu$ be such that $|\nu|=r$. In the same way as we did in (12), we have

$$
\begin{align*}
\left.\frac{1}{\omega(B)} \int_{B} \right\rvert\, S_{\alpha} f(x) & -S_{\alpha} f(\nu) \left\lvert\, d x \leq \frac{1}{\omega(B)} \int_{B} \int_{|y| \leq|x|} K_{\nu}(x, y) \frac{|f(y)|}{|y|^{n-\alpha}} d y d x\right.  \tag{19}\\
& +\frac{1}{\omega(B)} \int_{B} \int_{|x|<|y| \leq r} K_{\nu}(x, y) \frac{|f(y)|}{|y|^{n-\alpha}} d y d x
\end{align*}
$$

where $K_{\nu}(x, y)=\min \left\{1, \frac{|y|^{n-\alpha}}{|x|^{n-\alpha}}\right\}-\min \left\{1, \frac{| |^{n-\alpha}}{\left|| |^{n-\alpha}\right.}\right\}$.
We estimate the first term of (19). Let $B_{j}=B\left(0,2^{-j} r\right), j=0,1, \ldots$. Thus, since $K_{\nu}(x, y) \leq|y|^{n-\alpha}|x|^{-(n-\alpha)}$ for $|y| \leq|x|$ and (5), we have

$$
\begin{align*}
\frac{1}{\omega(B)} \int_{B} \int_{|y| \leq|x|} K_{\nu}(x, y) \frac{|f(y)|}{|y|^{n-\alpha}} d y d x & \leq \frac{1}{\omega(B)} \int_{B} \frac{1}{|x|^{n-\alpha}} \int_{|v| \leq|x|}|f(y)| d y d x \\
& \leq C \mid f \|_{B M_{0}^{y}(\omega)} \frac{1}{\omega(B)} \int_{B} \frac{\omega(B(0,|x|))}{|x|^{n-\alpha-n \gamma}} d x \\
& \leq C\|f\|_{B M_{0}^{\gamma}(\omega)} \frac{r^{n \gamma+\alpha}}{\omega(B)} \sum_{j=0}^{\infty} \frac{\omega\left(B_{j} \backslash B_{j+1}\right)}{2^{j(n \gamma+\alpha)}}  \tag{20}\\
& \leq C\|f\|_{B M_{0}^{\gamma}(\omega)} \frac{|B|^{\delta}}{\omega(B)} \sum_{j=0}^{\infty} \omega\left(B_{j} \backslash B_{j+1}\right) \\
& =C\|f\|_{B M_{0}^{\gamma}(\omega)}|B|^{\delta} .
\end{align*}
$$

For the second term of (19), since $0 \leq K_{\nu}(x, y) \leq 1$, we have

$$
\begin{align*}
\frac{1}{\omega(B)} \int_{B} \int_{|x| \leq|y| \leq r} K_{\nu}(x, y) \frac{|f(y)|}{|y|^{n-\alpha}} d y d x & \leq \frac{1}{\omega(B)} \int_{B} \frac{|f(y)|}{|y|^{n-\alpha}} \int_{|x| \leq|y|} 1 d x d y  \tag{21}\\
& \leq C \mid f \|_{\left.B M_{0}^{\gamma}(\omega)\right|^{|B|} .}
\end{align*}
$$

Therefore, by (19)-(21) we have proved

$$
\begin{equation*}
\frac{1}{\omega(B)|B|^{\delta}} \int_{B}\left|S_{\alpha} f(x)-S_{\alpha} f(\nu)\right| d x \leq C\|f\|_{B M_{0}^{\gamma}(\omega)}, \tag{22}
\end{equation*}
$$

for every ball $B$ centered at the origin.
We now consider $B=B\left(x_{0}, r\right)$ with $r<\left|x_{0}\right| / 8$. By (22) it is enough to consider only these balls $B$. Let $x \in B$ and $\nu=\left(\left|x_{0}\right|+r\right) x_{0} /\left|x_{0}\right|$. In the same way as we obtained (11) and (15) in the previous proof, we have

$$
S_{\alpha} f(x)-S_{\alpha} f(\nu)=\int_{|y| \leq|\nu|} K_{\nu}(x, y) \frac{f(y)}{|y|^{n-\alpha}} d y
$$

and $K_{\nu}(x, y) \leq \operatorname{Cr}|y|^{n-\alpha}|\nu|^{-(n-\alpha+1)}$. By $\omega \in D_{\eta}$, we have

$$
\begin{align*}
\frac{1}{\omega(B)} \int_{B}\left|S_{\alpha} f(x)-S_{\alpha} f(\nu)\right| d x & \leq C \frac{r^{n+1}}{\omega(B)|\nu|^{n-\alpha+1}} \int_{|\nu| \leq|\nu|}|f(y)| d y \\
& \leq C\|f\|_{B M_{0}^{\gamma}(\omega)} \frac{r^{n+1}}{|\nu|^{n-\alpha+1-n \gamma}} \frac{\omega(B(0,|\nu|))}{\omega(B)}  \tag{23}\\
& \leq C\|f\|_{B M_{0}^{\gamma}(\omega)}|B|^{\delta} .
\end{align*}
$$

Therefore, (22) and (23), complete the proof of the theorem.
Let $x, \nu \in \mathbb{R}^{n}, \nu \neq 0$, then

$$
\begin{align*}
\left|H_{\alpha} f(x)-H_{\alpha} f(\nu)\right| & \left.\leq \int_{|y| \leq|\nu|}|f(y)| \frac{1}{(|x|+|y|)^{n-\alpha}}-\frac{1}{(|\nu|+|y|)^{n-\alpha}} \right\rvert\, d y \\
& +\int_{|\nu|>|\nu|}|f(y)|\left|\frac{1}{(|x|+|y|)^{n-\alpha}}-\frac{1}{(|\nu|+|y|)^{n-\alpha}}\right| d y . \tag{24}
\end{align*}
$$

Proof of Theorem 1.6: We begin showing the sufficient condition. Let $B=$ $B\left(x_{0}, r\right)$ and let $u, v, U$ and $V$ as in (9) of the proof of Theorem 1.4. Note that if $x \in U$, $z \in V$ then for all $y \in \mathbb{R}^{n}$,

$$
\begin{equation*}
\frac{1}{(|x|+|y|)^{n-\alpha}}-\frac{1}{(|z|+|y|)^{n-\alpha}} \geq C \frac{r}{\left(\left|x_{0}\right|+r+|y|\right)^{n-\alpha+1}} . \tag{25}
\end{equation*}
$$

Hence, if $f$ is a nonnegative function in $L^{p}\left(\omega^{-p}\right)$ such that $\operatorname{supp}(f) \subset A(0, r, m)$ and taking $x_{0}=0$ in (25), we have

$$
\begin{equation*}
\left\|H_{\alpha} f\right\|_{B M O^{\delta}(\omega)} \geq \frac{C r^{n+1}}{\omega(B)|B|^{\delta}} \int_{A(0, r, m)} \frac{f(y)}{|y|^{n-\alpha+1}} d y \tag{26}
\end{equation*}
$$

for every ball $B$ centered at the origin.
Thus, taking $f_{m, j}(y)=|y|^{-(n-\alpha+1) /(p-1)} \omega^{p^{\prime}}(y) \chi_{A_{m, j}}(y)$ in (26) where $A_{m, j}=$ $A(0, r, m) \cap\{y: 1 / j \leq \omega(y)<j\}, m, j=1,2, \ldots$, using the boundedness of $H_{\alpha}$ and letting $m \rightarrow \infty, j \rightarrow \infty$ we obtain that $\omega \in \mathrm{H}_{0}(\alpha, p)$.

On the other hand, if $f$ is a nonnegative function in $L^{p}\left(\omega^{-p}\right)$ such that $\operatorname{supp}(f) \subset B\left(0,2\left(\left|x_{0}\right|+r\right)\right)$, then by (25)

$$
\begin{equation*}
\left\|H_{\alpha} f\right\|_{B M O^{\delta}(\omega)} \geq \frac{C r^{n+1}}{\omega(B)|B|^{\delta}\left(\left|x_{0}\right|+r\right)^{n-\alpha+1}} \int_{B\left(0,2\left(\left|x_{0}\right|+r\right)\right)} f(y) d y \tag{27}
\end{equation*}
$$

Thus, taking $f_{j}(y)=\omega^{p^{\prime}}(y) \chi_{A_{j}}(y)$ in (27) where $A_{j}=$ $B\left(0,2\left(\left|x_{0}\right|+r\right)\right) \cap\{y: 1 / j \leq \omega(y)<j\}, j=1,2, \ldots$, and using the boundedness of $H_{\alpha}$, we have

$$
\left(\int_{A_{j}} \omega^{p^{\prime}}(y) d y\right)^{1 / p^{\prime}} \leq C\left(\frac{\left|x_{0}\right|+r}{r}\right)^{n-\alpha+1} \frac{\omega(B)}{|B|^{1 / p}}
$$

Letting $j \rightarrow \infty$ and taking $x_{0}=0$ in the last expression, we can obtain that $\omega \in R H_{0}\left(p^{\prime}\right)$. Then, applying Hölder's inequality, we obtain $\omega \in D_{\eta}$.

Now, let us show the necessary condition. Let $f \in L^{p}\left(\omega^{-p}\right)$ such that $H_{\alpha} f(x)$ is finite for some $x \neq 0$ and let $\omega$ such that $\omega \in \mathrm{H}_{0}(\alpha, p) \cap R H_{0}\left(p^{\prime}\right) \cap D_{\eta}$. Hence $H_{\alpha} f \in L_{l o c}^{1}\left(\mathbb{R}^{n}\right)$ since $(i)$ of Lemma 4.7. We begin considering $B=B(0, r), x \in B$ and $x \neq 0$. Let $\nu$ be such that $|\nu|=r$. We estimate the two terms of (24). By (8), we have

$$
\begin{align*}
\left.\frac{1}{\omega(B)} \int_{B} \int_{|\nu| \leq|\nu|} \right\rvert\, \frac{f(y)}{(|x|+|y|)^{n-\alpha}} & \left.-\frac{f(y)}{(|\nu|+|y|)^{n-\alpha} \mid} \right\rvert\, d y d x \\
& \leq \frac{C}{\omega(B)} \int_{B} \int_{B} \frac{|f(y)|}{|x|^{n-\alpha}} d y d x  \tag{28}\\
& \leq C\|f\|_{L^{p}\left(\omega^{-p}\right)} \int_{B} \frac{1}{|x|^{n-\alpha+n / p}} d x=C\|f\|_{L^{p}\left(\omega^{-p}\right)}|B|^{\delta} .
\end{align*}
$$

To analyze the second term of (24), we use the mean value theorem, then

$$
\left|\frac{1}{(|x|+|y|)^{n-\alpha}}-\frac{1}{(|\nu|+|y|)^{n-\alpha}}\right| \leq C \frac{r}{|y|^{n-\alpha+1}} .
$$

Thus, by (i) of Lemma 4.6

$$
\begin{align*}
\frac{1}{\omega(B)} \int_{B} \int_{|y|>|\nu|} \left\lvert\, \frac{f(y)}{||x|+|y|)^{n-\alpha}}-\frac{f(y)}{(|\nu|+|y|)^{n-\alpha} \mid} d y d x\right. & \leq \frac{C r}{\omega(B)} \int_{B} \int_{B^{c}} \frac{|f(y)|}{|y|^{n-\alpha+1}} d y d x  \tag{29}\\
& \leq C\|f\|_{L^{p}\left(\omega^{-p}\right)|B|^{\delta}} .
\end{align*}
$$

Therefore, by (24)-(29), we have proved

$$
\begin{equation*}
\frac{1}{\omega(B)|B|^{\delta}} \int_{B}\left|H_{\alpha} f(x)-H_{\alpha} f(\nu)\right| d x \leq C\|f\|_{L^{p}\left(\omega^{-p}\right)}, \tag{30}
\end{equation*}
$$

for every ball $B$ centered at the origin.
We now consider $B=B\left(x_{0}, r\right)$ with $r<\left|x_{0}\right| / 8$. By (28) it is enough to consider only these balls $B$. Let $x \in B$ and $\nu=\left(\left|x_{0}\right|+r\right) x_{0} /\left|x_{0}\right|$, then $|\nu| \sim|x|$ and $|x| \sim\left|x_{0}\right|$. Using $|y| \leq|\nu|$ and the mean value theorem

$$
\left|\frac{1}{(|x|+|y|)^{n-\alpha}}-\frac{1}{(|\nu|+|y|)^{n-\alpha}}\right| \leq C \frac{r}{\left|x_{0}\right|^{n-\alpha+1}} .
$$

Then, by (8) and $\omega \in D_{\eta}$

$$
\begin{align*}
& \frac{1}{\omega(B)} \int_{B} \int_{|\nu| \leq|\nu|}|f(y)|\left|\frac{1}{(|x|+|y|)^{n-\alpha}}-\frac{1}{(|\nu|+|y|)^{n-\alpha}}\right| d y d x \\
& \leq \frac{C r^{n+1}}{\omega(B)\left|x_{0}\right|^{n-\alpha+1}} \int_{|\nu| \leq|\nu|}|f(y)| d y  \tag{31}\\
& \leq C\|f\|_{L^{p}\left(\omega^{-p}\right)} \frac{r^{n+1}}{\left|x_{0}\right|^{n-\alpha+1+n / p}} \omega(B(0,|\nu|)) \\
& =C\|f\|_{L^{p}\left(\omega^{-p}\right) \mid}|B|^{\delta} .
\end{align*}
$$

Now, using the mean value theorem

$$
\left|\frac{1}{(|x|+|y|)^{n-\alpha}}-\frac{1}{(|\nu|+|y|)^{n-\alpha}}\right| \leq C \frac{r}{|y|^{n-\alpha+1}} .
$$

Then, by (i) of Lemma 4.6

$$
\begin{align*}
\frac{1}{\omega(B)} \int_{B} \int_{|y|>|\nu|}\left|\frac{f(y)}{| | x|+|y|)^{n-\alpha}}-\frac{f(y)}{(|\nu|+|y|)^{n-\alpha} \mid}\right| d y d x & \leq \frac{C r}{\omega(B)} \int_{B} \int_{B^{c}} \frac{|f(y)|}{|y|^{n-\alpha+1}} d y d x  \tag{32}\\
& =C\|f\|_{L^{p}\left(\omega^{-p}\right)|B|^{\delta}} .
\end{align*}
$$

Therefore, by (24) with $\nu=\left(\left|x_{0}\right|+r\right) x_{0} /\left|x_{0}\right|$, (31) and (32), we have

$$
\frac{1}{\omega(B)|B|^{\delta}} \int_{B}\left|H_{\alpha} f(x)-H_{\alpha} f(\nu)\right| d x \leq C\|f\|_{L^{p}\left(\omega^{-p}\right)}
$$

for every ball $B=B\left(x_{0}, r\right)$ considered. This completes the proof of the theorem.
Proof of Theorem 1.7: We begin showing the sufficient condition. Let $B=B\left(x_{0}, r\right)$ and let $u, v, U$ and $V$ as in (9) of the proof of Theorem 1.4. Then, as in (26) of the proof of Theorem 4 (with $x_{0}=0$ ), we again have

$$
\begin{equation*}
\left\|H_{\alpha} f\right\|_{B M O^{\delta}(\omega)} \geq \frac{C r^{n+1}}{\omega(B)|B|^{\delta}} \int_{A(0, r, m)} \frac{f(y)}{|y|^{n-\alpha+1}} d y \tag{33}
\end{equation*}
$$

for every nonnegative function $f$ in $B M_{0}^{\gamma}(\omega)$ such that $\operatorname{supp}(f) \subset A(0, r, m)$ and for every ball $B$ centered at the origin.

Thus, taking $f(y)=|y|^{n \gamma} \omega(y) \chi_{A(0, r, m)}(y)$ in (33), using that $\|f\|_{B M_{0}^{\gamma}(\omega)} \leq 1$, the boundedness of $H_{\alpha}$ and letting $m \rightarrow \infty$, we have that $\omega \in \mathrm{H}_{0}(\alpha+n \gamma, \infty)$.

On the other hand, as in (27) of the proof of Theorem 1.6 we again have

$$
\begin{equation*}
\left\|H_{a} f\right\|_{B M O^{\delta}(\omega)} \geq \frac{C r^{n+1}}{\omega(B)|B|^{\delta}\left(\left|x_{0}\right|+r\right)^{n-\alpha+1}} \int_{B\left(0,2\left(\left|x_{0}\right|+r\right)\right)} f(y) d y \tag{34}
\end{equation*}
$$

for every nonnegative function $f$ in $B M_{0}^{\gamma}(\omega)$ such that $\operatorname{supp}(f) \subset B\left(0,2\left(\left|x_{0}\right|+r\right)\right)$ and for every ball $B=B\left(x_{0}, r\right)$.

If $\gamma=0$, we take $f(y)=\omega(y) \chi_{B\left(0,2\left(\left|x_{0}\right|+r\right)\right)}(y)$ in (4.34) and since $\|f\|_{B M_{0}^{\gamma}(\omega)} \leq 1$ and the boundedness of $H_{\alpha}$, we have $\omega \in D_{\eta}$.

If $\gamma>0$, let $f(y)=P_{n \gamma}\left(\omega \chi_{B\left(0,2\left(\left|x_{0}\right|+r\right)\right)}\right)(y)$ then $\|f\|_{B M_{0}^{\gamma}(\omega)} \leq C$ and as in (4.18) of the proof of Theorem 1.5, we have

$$
\int_{B\left(0,2\left(\left|x_{0}\right|+r\right)\right)} f(y) d y \geq C\left(\left|x_{0}\right|+r\right)^{n \gamma} \omega\left(B\left(0,\left|x_{0}\right|+r\right)\right) .
$$

Therefore, using this function $f$ in (34) and the boundedness of $H_{\alpha}$, we have $\omega \in D_{\eta}$.
Now, let us show the necessary condition. Let $f \in B M_{0}^{\gamma}(\omega)$ such that $H_{\alpha} f(x)$ is finite for some $x \neq 0$ and let $\omega \in \mathrm{H}_{0}(\alpha+n \gamma, \infty) \cap D_{\eta}$. Hence $H_{\alpha} f \in L_{l o c}^{1}\left(\mathbb{R}^{n}\right)$ by (ii) of Lemma 4.7. We begin considering $B=B(0, r), x \in B$ and $x \neq 0$. Let $\nu$ be such that $|\nu|=r$. We estimate the two terms of (24). Then,

$$
\begin{align*}
& \frac{1}{\omega(B)} \int_{B} \int_{|v| \leq|\nu|}|f(y)|\left|\frac{1}{(|x|+|y|)^{n-\alpha}}-\frac{1}{(|\nu|+|y|)^{n-\alpha}}\right| d y d x \\
& \leq \frac{C}{\omega(B)} \int_{B} \frac{1}{|x|^{n-\alpha}} \int_{|y| \leq|\nu|} f(y) d y d x  \tag{35}\\
& \leq C\|f\|_{B M_{0}^{\gamma}(\omega)} \frac{1}{\omega(B)} \int_{B} \frac{\omega(B(0,|\nu|))|\nu|^{n \gamma}}{|x|^{n-\alpha}} d x \\
& \leq C\|f\|_{B M_{0}^{\gamma}(\omega)}|B|^{\delta} .
\end{align*}
$$

For the second term of (24), using the mean value theorem

$$
\begin{equation*}
\left|\frac{1}{(|x|+|y|)^{n-\alpha}}-\frac{1}{(|\nu|+|y|)^{n-\alpha}}\right| \leq C \frac{r}{|y|^{n-\alpha+1}} . \tag{36}
\end{equation*}
$$

Then, by (ii) of Lemma 4.6

$$
\begin{align*}
& \frac{1}{\omega(B)} \int_{B} \int_{|y|>|\nu|}|f(y)| \left\lvert\, \frac{1}{(|x|+|y|)^{n-\alpha}}-\frac{1}{(|\nu|+|y|)^{n-\alpha} \mid} d y d x\right. \\
& \leq \frac{C r}{\omega(B)} \int_{B} \int_{B^{c}} \frac{|f(y)|}{|y|^{n-\alpha+1}} d y d x  \tag{37}\\
& \leq C\|f\|_{B M_{0}^{\gamma}(\omega)}|B|^{\delta} .
\end{align*}
$$

Therefore, by (24) and (35)-(37), we have proved

$$
\begin{equation*}
\frac{1}{\omega(B)|B|^{\delta}} \int_{B}\left|H_{\alpha} f(x)-H_{\alpha} f(\nu)\right| d x \leq C\|f\|_{B M_{0}^{\%}(\omega)}, \tag{38}
\end{equation*}
$$

for every ball $B$ centered at the origin.

We now consider $B=B\left(x_{0}, r\right)$ with $r<\left|x_{0}\right| / 8$. By (33) it is enough to consider only these balls $B$. Let $x \in B$ and $\nu=\left(\left|x_{0}\right|+r\right) x_{0} /\left|x_{0}\right|$, then $|\nu| \sim|x|$ and $|x| \sim\left|x_{0}\right|$. If $|y| \leq|\nu|$, by the mean value theorem

$$
\left|\frac{1}{(|x|+|y|)^{n-\alpha}}-\frac{1}{(|\nu|+|y|)^{n-\alpha}}\right| \leq C \frac{r}{\left|x_{0}\right|^{n-\alpha+1}} .
$$

Then, since $\omega \in D_{\eta}$, we have

$$
\begin{align*}
& \frac{1}{\omega(B)} \int_{B} \int_{|y| \leq|\nu|}|f(y)| \left\lvert\, \frac{1}{(|x|+|y|)^{n-\alpha}}-\frac{1}{(|\nu|+|y|)^{n-\alpha} \mid} d y d x\right. \\
& \leq \frac{C r^{n+1}}{\omega(B)\left|x_{0}\right|^{n-\alpha+1}} \int_{|y| \leq|\nu|}|f(y)| d y  \tag{39}\\
& \leq\|f\|_{B M_{0}^{*}(\omega)}|B|^{\delta} .
\end{align*}
$$

On the other hand, using again the mean value theorem as in (36) and (ii) of Lemma 4.6, we get

$$
\begin{align*}
& \frac{1}{\omega(B)} \int_{B} \int_{|y|>|\nu|}|f(y)|\left|\frac{1}{(|x|+|y|)^{n-\alpha}}-\frac{1}{(|\nu|+|y|)^{n-\alpha}}\right| d y d x \\
& \leq \frac{C r^{n+1}}{\omega(B)} \int_{|\nu|>|\nu|} \frac{|f(y)|}{|y|^{n-\alpha+1}} d y  \tag{40}\\
& \leq C\|f\|_{B M_{0}^{\%}(\omega)} \frac{r^{n+1}}{\omega(B)} \frac{\omega(B(0,|\nu|))}{|\nu|^{n \eta}} \\
& \leq C\|f\|_{B M_{0}^{\gamma}(\omega)}|B|^{\delta} .
\end{align*}
$$

Thus, by (24) and (39)-(40), we have proved

$$
\frac{1}{\omega(B)|B|^{\delta}} \int_{B}\left|H_{\alpha} f(x)-H_{\alpha} f(\nu)\right| d x \leq C\|f\|_{B M_{0}^{\gamma}(\omega)},
$$

for every ball $B=B\left(x_{0}, r\right)$ considered. This completes the proof of the theorem.

## 5. Conclusions

As a conclusion to this chapter, we have given necessary and sufficient conditions for the generalized Calderón and Hilbert operators to be bounded from weighted Lesbesgue spaces into suitable weighted $B M O$ and Lipschitz spaces. Then, we have obtained results on the boundedness of these operators from $L^{\infty}$ into $B M O$, even in the unweighted case for the Hilbert operator. The class of weights involved are close to the doubling and reverse Hölder conditions related to the Muckenhoupt's classes.

The study of the weighted boundedness for integral operators on function spaces, like the one we develop in this chapter, is one of the main research fields in harmonic analysis. In particular, it has had a profound influence in partial differential equations, several complex variables, and number theory. Evidence of such success and
importance is the pioneering work of leading mathematicians Bourgain, Zygmund, Calderón, Muckenhoupt, Wheeden, C. Fefferman, Stein, Ricci, Tao and so on.

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Section 2
Applications

## Chapter 4

# Effect of Titanium Oxide Nanofluid over Cattaneo-Christov Model 

Hammad Khalil, Tehseen Zahra, Zaffer Elahi and Azeem Shahzad


#### Abstract

The proposed chapter deals with the study of heat transfer development of titanium oxide nanofluid of platelet shape nanoparticles over a vertical stretching cylinder. The set of nonlinear equations is obtained using suitable transformation on the governing equations that are then solved with numerical scheme BVP4C. The obtained results are interpreted graphically and numerically. The effects of Prandtl, Eckert, and unsteadiness parameters on temperature distribution are depicted. Moreover the skin friction and Nusselt number are also computed.


Keywords: Cattaneo-Christov model, heat transfer, vertical cylinder

## 1. Introduction

Cattaneo-Christov model is an improved version of Fourier law as Fourier law does not detect the initial temperature disturbance; to overcome this ambiguity, Cattaneo added a thermal relaxation parameter. This parameter covers the ambiguity of Fourier law. The classical Fourier law is obtained while vanishing the relaxation parameter [1]. Cattaneo-Christov heat flux model gives us heat transfer rate in stretching cylinders as well as sheets. Heat transfer is a wonderful natural phenomenon that occurs when two bodies have a thermal difference until both bodies are at thermal equilibrium. The Cattaneo-Christov model is in the form of a heat equation. The thermal convection effect is studied using the Christov heat model in conjunction with the Cattaneo heat model [2]. It has been realized that the development of stretchy surfaces and the flow field that surrounds them speaks to a variety of technological and industrial applications, such as paper making, glass blowing, crystal growth, and aerodynamic plastic sheet extrusion [3]. Heat transfer is a common natural occurrence as long as there is a temperature differential between things or between various regions of the same object, heat transfer will occur. As a result, a lot of effort has gone into predicting the heat transport behavior. In several starting and boundary problems, the uniqueness and structural stability of the solutions for the temperature governing equations using the Cattaneo-Christov heat flow model have been demonstrated. The chapter released uses the Cattaneo-Christov heat flux model to analyze the flow and heat transfer of upper-convective Maxwell fluid across a stretching sheet [4]. Efforts have been undertaken to increase the thermal efficiency of processes during the last many decades. On the one hand, there has been an
attempt to lower the size of the equipment by increasing the thermal exchange surface, such as with fins, and on the other hand, novel fluid exchangers with higher thermal conductivity have been developed. Different NPS types (metallic, nonmetallic, and carbon based) have been synthesized and dispersed in conventional fluids such as water, oil, or ethylene glycol referred to as nanofluids since the advent of nanotechnology and the possibility of synthesizing materials on a nanometric scale [5]. The boundary layer flow and heat transfer caused by stretching flat plates or cylinders are both practical and theoretically interesting in fiber technology and extrusion operations. This method is used to produce polymer sheets and plastic films. The cooling of an infinite metallic plate in a cooling bath, the boundary layer along material handling conveyors, the aerodynamic extrusion of plastic sheets, the boundary layer along a liquid film in condensation processes, paper production, glass blowing, metal spinning and drawing plastic films, and polymer extrusion are all examples of boundary layers [6]. The aim of this chapter is to manipulate the heat transfer rate of titanium oxide nanofluid with the CattaneoChristov heat flux model over a vertical stretching cylinder.

## 2. Mathematical formulation

In the coordinate plane, assume that the cylinder is taken in the vertical direction along the $z$-axis, and the $r$-axis is normal to the axis of the cylinder. Consider the fluid is moving with surface velocity,

$$
U_{w}=\frac{b z}{1-\alpha t}
$$

In the direction of stretching cylinder under the external magnetic field defined by

$$
B(t)=\frac{B_{0}}{\sqrt{1-\alpha t}}
$$

Further, suppose $u=u(r, z, t)$ and $w=w(r, z, t)$ be the velocity components along the respective axes of the coordinate plane, while $T=T(r, z, t)$ denotes the temperature of the nanofluid, as shown in Figure 1.


Figure 1.
Representation of the dilemma.

Let the base fluid (water) and platelet-shaped NPS in thermal equilibrium. Under these assumptions, the equation of continuity, momentum equation, and energy equations are obtained, which are as follows:

$$
\begin{gather*}
\frac{\partial(r u)}{\partial r}+\frac{\partial(r w)}{\partial z}=0  \tag{1}\\
+\frac{u \partial w}{\partial r}+\frac{w \partial w}{\partial z}=\frac{\nu}{r} \frac{\partial}{\partial r}\left(\frac{r \partial w}{\partial r}\right)-\frac{\sigma}{\rho} B_{0}^{2} w+\mathrm{g} \beta(T-T \infty)  \tag{2}\\
\frac{\partial T}{\partial t}+\frac{u \partial T}{\partial r}+\frac{w \partial T}{\partial z}= \\
\frac{\nu}{C p}\left(\frac{\partial w}{\partial r}\right)^{2}+\lambda_{1}\left[\frac{2 \nu}{C p}\left(\frac{\partial w}{\partial r}\right)\left(\frac{\partial^{2} w}{\partial t \partial r}+u \frac{\partial^{2} w}{\partial r^{2}}\right)-\left(\frac{\partial^{2} T}{\partial t^{2}}\right.\right. \\
 \tag{3}\\
+u^{2} \frac{\partial^{2} T}{\partial r^{2}}+w^{2} \frac{\partial^{2} T}{\partial z^{2}} 2 u \frac{\partial^{2} T}{\partial t \partial r}+2 w \frac{\partial^{2} T}{\partial t \partial z}+2 u w \frac{\partial^{2} T}{\partial z \partial r} \\
\\
\left.\left.+\left(\frac{\partial u}{\partial t}+u \frac{\partial u}{\partial r}+w \frac{\partial u}{\partial z}\right) \frac{\partial T}{\partial r}+\left(\frac{\partial w}{\partial t}+u \frac{\partial w}{\partial r}+w \frac{\partial w}{\partial z}\right) \frac{\partial T}{\partial z}\right)\right]
\end{gather*}
$$

Subject to the boundary conditions

$$
\begin{gather*}
\mathrm{u}=0, \mathrm{w}=U_{\mathrm{w}}, \frac{\partial T}{\partial r}=0 \text { at } \mathrm{r}=\mathrm{R}, \\
w \rightarrow o, T \rightarrow T_{\infty} \text { as } r \rightarrow \infty \tag{4}
\end{gather*}
$$

The thermophysical properties of density ( $\rho n f$ ), dynamic viscosity ( $\mu n f$ ), electric conductivity ( $\sigma n f$ ), diffusivity ( $\alpha n f$ ), and heat capacity ( $\rho C p$ ) can be defined in Refs. $[7,8]$, while the ratio of thermal conductivity of nanofluid and base fluid is given by the following equation:

$$
\begin{equation*}
\frac{k_{n f}}{k_{f}}=\left[\frac{k_{s}+(m-1) k_{f}+(m-1)\left(k_{s}-k_{f}\right) \phi}{k_{s}+(m-1) k_{f}-\left(k_{s}-k_{f}\right) \phi}\right] \tag{5}
\end{equation*}
$$

where $\phi$ denotes volume-fraction of NPS.
The thermophysical properties of titanium nanofluid with base fluid as water [9] are given in Table 1, while the viscosity coefficients $A 1, A 2$, and shape factor $m$ values of $\mathrm{TiO}_{2}$ nanofluid [10] are listed in Table 2.

| Base | Density $\left(\mathbf{k g} / \mathbf{m}^{\mathbf{3}}\right)$ | Thermal conductivity (W/m K) | Specific heat $(\mathbf{J} / \mathbf{k g}$ K) |
| :---: | :---: | :---: | :---: |
| $\mathrm{TiO}_{2}$ | 3900 | 8.4 | 0.8692 |
| $\mathrm{H}_{2} \mathrm{O}$ | 997.1 | 0.613 | 4179 |

Table 1.
Thermophysical properties of base fluid and $\mathrm{TiO}_{2}$ nanoparticles.

|  | Platelet |  |
| :--- | :---: | :---: |
| $\mathrm{A}_{1}$ | $\mathrm{~A}_{2}$ | m |
| 37.1 | 612.6 | 5.72 |

Table 2.
Viscosity and shape factor values of platelet-shaped nanoparticles.

Introducing the transformations, as

$$
\begin{equation*}
\mathrm{T}=\mathrm{T}_{\infty}+\left(\mathrm{Tw}-\mathrm{T}_{\infty}\right) \theta(\eta), \eta=\left(\sqrt{\frac{c}{\nu(1-\alpha t)}}\right)\left(\frac{r^{2}-R^{2}}{2 R}\right), \psi=\left(\sqrt{\frac{c \nu}{(1-\alpha t)}}\right) z v f(\eta) \partial \psi \tag{6}
\end{equation*}
$$

where $\psi$ is the stream function (describes the flow pattern) and is defined as $u=$ $\frac{-1}{r} \frac{\partial \psi}{\partial r}$ and $\mathrm{w}=\frac{1}{r} \frac{\partial \psi}{\partial r}$. The governing Eqs. (2)-(5) have been transformed to Eqs. (8)-(10) using similarity variables in Eq. (7), as

$$
\begin{gather*}
\epsilon_{1}(1+2 C \eta) f^{\prime \prime \prime}(\eta)+2 \epsilon_{1} C f^{\prime \prime}(\eta)-\epsilon_{3} M f(\eta) \\
+\left[f(\eta) f^{\prime \prime \prime}(\eta)-f^{\prime 2}(\eta)-S\left(f^{\prime}(\eta)+\frac{\eta}{2} f^{\prime \prime}(\eta)\right)\right]+\lambda \theta(\eta)=0  \tag{7}\\
(1+2 C \eta)\left(\epsilon_{1} E c f^{\prime \prime 2}(\eta)+\frac{\epsilon_{2}}{\operatorname{Pr}} \theta^{\prime \prime}(\eta)\right)+\frac{2 \epsilon_{2}}{\operatorname{Pr}} C \theta^{\prime(\eta)}+f(\eta) \theta^{\prime(\eta)}-S\left(2 \theta(\eta)+\frac{\eta}{2} f^{\prime \prime}(\eta)\right) \\
+\beta\left[\begin{array}{c}
\epsilon_{1} \operatorname{Ec}(1+2 \eta)\left(3 S f^{\prime \prime 2}(\eta)+S \eta f^{\prime \prime}(\eta) f^{\prime \prime \prime}(\eta)-2 f(\eta) f^{\prime \prime}(\eta) f^{\prime \prime \prime}(\eta)\right) 2 \epsilon_{1} \operatorname{EcCf}(\eta) f^{\prime \prime 2}(\eta) \\
\left.-\left(\begin{array}{l}
S^{2}\binom{5 f^{\prime}(\eta) \theta(\eta)-\frac{11}{2} f(\eta) \theta^{\prime}(\eta)-\eta f(\eta) \theta^{\prime \prime}(\eta)}{6 \theta(\eta)+\frac{11}{4} \theta^{\prime \prime}(\eta) f(\eta)+S\left(\eta-\frac{1}{2}\right) f^{\prime}(\eta) \theta^{\prime \prime}(\eta)+\frac{\eta}{2} f^{\prime \prime}(\eta) \theta(\eta)} \\
+(\eta) \\
+f^{\prime}(\eta) \theta^{\prime \prime}(\eta)-f(\eta) f^{\prime}(\eta) \theta^{\prime}(\eta)-f(\eta) f^{\prime \prime}(\eta) \theta(\eta)+f^{\prime 2}(\eta) \theta(\eta)
\end{array}\right)\right]
\end{array}\right] \\
-(\eta) . \tag{8}
\end{gather*}
$$

where $\beta$ is the thermal relaxation parameter and is given by,

$$
\begin{equation*}
\beta 1=\frac{c \lambda 1}{(1-\alpha t)} \tag{9}
\end{equation*}
$$

Under the boundary condition,

$$
\begin{gather*}
f(0)=0, f^{\prime}(0)=1, \theta^{\prime}(0)=\frac{-k_{f}}{k_{n f}} \gamma(1-\theta(0)) \text { at } \eta=R \\
f^{\prime}(\eta)=0, \theta(0) \text { as } \eta \rightarrow \infty \tag{10}
\end{gather*}
$$

Now the dimensionless constants, such as $E c, \operatorname{Pr}, \phi, M$, and $S$, and that of $\epsilon 1, \epsilon 2$, and $\epsilon 3$ are used frequently in the above equations, defined in Ref. [11]. For various
values of dimensionless parameters, the value of the local Nusselt number is shown in Table 3. Nusselt number can be defined as,

$$
\begin{equation*}
\mathrm{Nu}=\frac{z k_{n f}}{k_{f}\left(T_{w-} T_{\infty}\right)}\left[\frac{\partial T}{\partial r}\right]_{\mathrm{r}=\mathrm{R}} \tag{11}
\end{equation*}
$$

The non-dimensionless form of Eq. (11), using Eq. (6), as

$$
\begin{equation*}
R e^{\frac{-1}{2}} N u=-\frac{k_{n f}}{k_{f}} \theta^{\prime}(0) \tag{12}
\end{equation*}
$$

## 3. Method of solution

To find the numerical solution of a nonlinear system (7) and (8), the set of firstorder linear equations is obtained by considering the following assumptions. By putting these assumptions in the above equations, we get first-order linear equations, which are then used in MATLAB by using BVP4C scheme to get numerical and graphical results.

$$
\begin{gather*}
y_{1}=f,  \tag{13}\\
y_{1}^{\prime}=y_{2},  \tag{14}\\
y_{2}^{\prime}=y_{3},  \tag{15}\\
y_{3}^{\prime}=g_{1},  \tag{16}\\
\theta=y_{4},  \tag{17}\\
y_{4}^{\prime}=g_{1},  \tag{18}\\
y_{5}^{\prime}=g_{2},  \tag{19}\\
y_{1}(0)=0, y_{2}(0)=1, y_{4}(0)-1=0 \text { at } \eta=0  \tag{20}\\
y_{2}(\eta)=0, y_{4}(\eta)=0 . \text { as } \eta \rightarrow \infty \tag{21}
\end{gather*}
$$

|  | Physical parameters | Platelet |
| :--- | :---: | :---: |
| $\boldsymbol{E c}$ | $\operatorname{Pr}$ | $\mathbf{R e}^{\frac{-1}{2}} \mathbf{N u}$ |
| 0.0 | 6.0 | 0.037666741 |
| 0.5 | - | 0.038924615 |
| 1.0 | - | 0.039829280 |
| 1.0 | 4.0 | 0.042765482 |
| - | 6.0 | 0.038924615 |
| - | 8.0 | 0.036302577 |

Table 3.
Nusselt number of platelet shape nanoparticle.
where

$$
\begin{equation*}
g_{1}=\frac{1}{\epsilon_{1}(1+2 \eta C)}\left(\epsilon_{3} M y_{2}-2 \epsilon_{1} C y_{3}-\lambda y_{4}-\left(y_{1} y_{3}-y_{2}^{2}-S\left(y_{2}+\frac{\eta}{2} y_{3}\right)\right)\right) \tag{22}
\end{equation*}
$$

and

$$
\begin{align*}
& g_{2}=\frac{1}{\frac{(1+2 \eta C) \epsilon_{3}}{P r}-\beta_{1}\left(\frac{(S \eta)^{2}}{4}-S \eta y_{1}+y_{1}^{2}\right)}\left(S\left(y_{4}+\frac{\eta}{2} y_{5}\right)-y_{1}+y_{2} y_{4}-\left(\frac{2 \epsilon_{2} C y_{5}}{P r}\right)\right. \\
& -(1+2 \eta C) \epsilon_{1} E c y_{3}^{2}-\beta_{1}\left(\epsilon_{1} E c(1+2 \eta C)\left(3 S y_{3}^{2}\right)+S t y_{3} y_{3}^{\prime}-2 \epsilon_{1} E c C y_{1} y_{3}^{2}-2 y_{1} y_{3} y_{3}^{\prime}\right) \\
& \left.\left.-\left(S^{2}\left(6 y_{4}+\frac{11}{4} \eta y_{5}\right)-S\left(5 y_{2} y_{4}-\frac{11}{2} y_{1} y_{5}\right)+\left(\eta-\frac{1}{2} y_{2} y_{5}+\frac{\eta}{2} y_{3} y_{4}\right)+y_{1} y_{3} y_{5}+y_{2}^{2} y_{4}\right)\right)\right) . \tag{23}
\end{align*}
$$

$g_{1}$ and $g_{2}$ are the obtained first-order linear equations.

## 4. Analysis of results

Obtained numerical results and the effect of various parameters on temperature profile are obtained in both numerical and graphical form and discussed in detail.

### 4.1 Graphical analysis

Figure 2 describes the influence of Eckert number for platelet shape nanoparticle. By varying the value of Eckert, it can be seen that the temperature is increasing. Physically, it can be seen that the Eckert number enhances the thermal conductivity of the fluid. Figure 3 shows that the increase in the value of the Prandtl number results


Figure 2.
Effect of Eckert number on the temperature profile.


Figure 3.
Effect of Prandtl number on the temperature profile.


Figure 4.
Effect of unsteadiness parameter number on the temperature profile.
in deceleration in temperature because of the reduction in thermal diffusivity. As the unsteadiness parameter increases, the temperature gradually decreases, as shown in Figure 4. Physically the value of the unsteadiness parameter is grown up, and the thickness of the thermal boundary layer decreases, which results in a decline in the temperature profile.

### 4.2 Numerical results

The heat transfer rate is calculated for platelet shape nanoparticles, which is given in Table 3. It is inferred that with the rise in Eckert number, Nusselt number increases while the reverse trend is seen for Prandtl number

## 5. Conclusion

By numerical computation, the effect of platelet shape nanoparticle on $\mathrm{TiO}_{2}$ nanofluid over a vertical stretching cylinder is seen in this chapter. Influence of different physical parameters, such as Eckerd and Prandtl numbers on temperature profile, is examined both graphically and numerically. The Nusselt number increases for Eckert number Ec, which decreases for Prandtl number Pr. Graphical result shows acceleration in temperature profile, while the reverse trend is found in Figure 2.

## Nomenclature

$u, w \quad$ velocity components along $r, z$ directions ( $\mathrm{m} / \mathrm{s}$ )
$\alpha f, \alpha n f \quad$ thermal diffusion of base fluid and nanofluid ( $\mathrm{m}^{2} / \mathrm{s}$ )
$\rho f, \rho n f \quad$ density of base fluid and nanofluid $\left(\mathrm{kg} / \mathrm{m}^{3}\right)$
$\mu f, \mu n f \quad$ viscosity of base fluid and nanofluid ( $\mathrm{kg} \mathrm{m} / \mathrm{s}$ )
$\nu f$, vnf kinematic viscosity of base fluid and nanofluid ( $\mathrm{m}^{2} / \mathrm{s}$ )
onf electrical conductivity (-)
$U w \quad$ surface velocity ( $\mathrm{m} / \mathrm{s}$ )
onf electrical conductivity (-)
$(\rho C p) n f \quad$ heat capacity of nanofluid (-)

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Chapter 5

# Nonlinear Dynamics Phenomenon in a Polydyne Cam with an Offset Flat Faced Follower Mechanism with Clearance 

Louay S. Yousuf


#### Abstract

Nonlinear response of the follower motion is simulated at different cam speeds, different coefficient of restitution, and different internal distance of the follower guide from inside. The nonlinear response of the follower is employed to investigate the chaotic phenomenon in cam follower system in the presence of follower offset. The numerical results are done using SolidWorks software. The chaos phenomenon is detected using Poincare' maps with phase-plane portraits, the largest Lyapunov exponent parameter, and bifurcation diagram. The largest Lyapunov exponent has a maximum values when the follower offsets to the right, while the largest Lyapunov exponent has a minimum values when the follower offsets to the left. The chaotic phenomenon in cam follower system when the follower offsets to the left is more than the chaotic phenomenon when the follower offsets to the right.


Keywords: chaotic phenomenon, follower offset, Lyapunov exponent parameter, nonlinear response. Poincare' maps

## 1. Introduction

The proposed system can be found in windshield wiper on the front window of the car in which the rotary motion of the cam transforms into an oscillating motion. Yang et al. introduced the mathematical model to describe the separation, transient impact, and contact in cam follower system using oblique impact, [1]. They showed that the cam and the follower system kept permanent contact without the use of coefficient of restitution at low speeds for the cam. Yousuf studied the detachment between the cam and the follower using largest Lyapunov exponent parameter, power density function of Fast Fourier Transform (FFT), and Poincare' maps due to the nonlinear dynamics phenomenon of the follower. Nonlinear response of the follower displacement is calculated at different cam speeds, different coefficient of restitution, different contact conditions, and different internal distance of the follower guide from inside [2, 3]. Flores et al. used a nonsmooth dynamics approach to model the interaction of the colliding bodies using Coulomb's law for dry friction [4]. Lassaad et al. studied the effect of cam profile error on the nonlinear dynamics behavior of oscillating roller
follower system by using a model with eight degrees of freedom of two nonlinear Hertzian contacts [5]. Li and Du used the coefficient of restitution as a main control parameter to analyze the periodic movement and the bifurcation region in Non-fixed constrained collision vibration system [6]. Wu et al. studied the influence of the joint clearance on the dynamic response of a planar mechanism with two driving links and prismatic pair clearance under variable input speeds [7]. They concluded that the largest Lyapunov exponents are dependent on the clearance size and the input speed. Chen et al. identified the chaos phenomenon of the 2-DOF nine-bar mechanism with a revolute clearance using the phase diagrams, the Poincaré portraits, and largest Lyapunov exponent parameter [8]. Bifurcation diagrams with changing clearance value, friction coefficient, and driving speed are drawn. The aim of this paper is to discuss the chaotic phenomenon of an offset follower through the use of impact coefficient of restitution at different follower guides' clearances and different cam speeds.

## 2. Numerical simulation

Follower displacement is calculated using SolidWorks software [9]. The follower moved with three degrees of freedom. Four values of the follower guide's from inside (I.D. $=16,17,18,19 \mathrm{~mm}$ ) at different cam speeds are used. The follower with the offset ( $\mathrm{O}=20,30,40,50 \mathrm{~mm}$ ) are chosen. The impact coefficient of restitution with the values ( $0.2,0.3$, and 0.4 ) is considered in the calculation of nonlinear response of the follower in the presence of follower offset. Cam follower mechanism is shown in Figure 1.


Figure 1.
Polydyne cam with an offset flat-faced follower.

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Figure 2.
Nonlinear response mapping when the follower offsets to the right ( $O=10 \mathrm{~mm}$ ).

The chaotic phenomenon in cam follower system is increased with the increasing of impact coefficient of restitution in which the impact will happen due the loss in potential energy of the follower and due to the increase in follower guide clearance value. Figures 2 and 3 show the mapping of nonlinear response of the follower at


Figure 3.
Nonlinear response mapping when the follower offsets to the left $(O=10 \mathrm{~mm})$.
different cam speeds, different follower guides' clearances, and different impact coefficient of restitution when the follower offsets to the right and left respectively ( $\mathrm{O}=10 \mathrm{~mm}$ ). The nonlinear response of the follower is periodic as shown in Figure 2a and both the cam and the follower are in permanent contact. The follower lost the contact with the cam at time ( $\mathrm{t}=13.58 \mathrm{~s}$ ) and $(\mathrm{t}=15.99 \mathrm{~s})$ at detachment height
( 26.98 mm ) and ( 27.43 mm ) respectively. Due to the coefficient of restitution, the follower keep bouncing from the cam from ( $t=0.36 \mathrm{~s}$ ) to ( $\mathrm{t}=5.658 \mathrm{~s}$ ) while the follower will regain energy and keep permanent contact with the cam for the period from ( $\mathrm{t}=9.208 \mathrm{~s}$ ) to ( $\mathrm{t}=10.11 \mathrm{~s}$ ) which is having a periodic motion as illustrated in Figure 2b. The chaotic motion is shown in Figure 2c-f which increased with the increasing of follower guides' clearances, cam speeds, and coefficient of restitution. There is an intangible impact when the coefficient of restitution (0.2) and the dissipation in potential energy is occurred due to sliding while the contact is still valid between the cam and the follower, as shown in Figure 3a. The periodic and chaotic motion is together shown in Figure $\mathbf{3 b}$ and $\mathbf{c}$. The periodic motion is shown from the period ( $\mathrm{t}=6.1 \mathrm{~s}$ ) to $(\mathrm{t}=10.26 \mathrm{~s})$ and from the period $(\mathrm{t}=14.14 \mathrm{~s})$ to $(\mathrm{t}=19.55 \mathrm{~s})$ as shown in Figure 3b while the periodic motion begins from the period ( $\mathrm{t}=1.264 \mathrm{~s}$ ) to ( $\mathrm{t}=3.808 \mathrm{~s}$ ) as shown in Figure 3c. The chaotic motion is shown in Figure 3d-f.


Figure 4.
Bifurcation diagram against cam speeds.


Figure 5.
Bifurcation diagram when the follower offsets to the left $(O=50 \mathrm{~mm})$.

## 3. Bifurcation diagram

The contrast in angular displacement for the cam and the follower is used in the calculation of bifurcation diagram [10, 11]. Figure 4 is built at the follower guide's from inside (I.D. $=19 \mathrm{~mm}$ ) when the follower offsets to the right and left $(\mathrm{O}=50 \mathrm{~mm})$.

The periodic motion is shown in Figure 4 in which it has the blue trend at cam speeds ( $\mathrm{N}=100-300 \mathrm{rpm}$ ) while the quasi-periodic motion of the follower has red trend at cam speeds ( $\mathrm{N}=100 \mathrm{rpm}$ ). The transition to chaos for the system when the follower offsets to the left is grown faster than the system when the follower offsets to the right as indicated in Figures 5 and 6. It can be concluded that the transition to chaos is incremented with the increment in cam speeds.

## 4. Lyapunov exponent parameter

Local Lyapunov exponent parameter is used to detect the chaotic phenomenon of nonlinear response of the follower attractor. Positive Lyapunov exponent refers to chaotic phenomenon while negative Lyapunov exponent indicates to periodic motion [12]. Figure 7 shows the local Lyapunov exponent against number of points when the follower offsets to the right ( $\mathrm{O}=10 \mathrm{~mm}$ ) at coefficient of restitution ( 0.2 ), cam speed ( $\mathrm{N}=200 \mathrm{rpm}$ ), and follower guide's clearance ( 16 mm ). In this figure there are positive and negative local Lyapunov exponent in which the negative local Lyapunov exponent represents to the steady state while the positive local Lyapunov exponent reflects to the transient state. Each value of local Lyapunov exponent has a value of embedding dimension [13].

## 5. Poincare' maps with phase-plane portraits

The contact status of the follower is detected using Poincare' map at high and low speeds [14]. Moreover, the quantity of the black dots in Poincare' maps detects the


Figure 6.
Bifurcation diagram when the follower offsets to the right ( $O=50 \mathrm{~mm}$ ).

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Figure 7.
Local Lyapunov exponent when the follower offsets to the right $(O=10 \mathrm{~mm})$ at cam speed $(N=200 \mathrm{rpm})$, coefficient of restitution ( 0.2 ), and follower guide's clearance (16 mm).


Figure 8.
Phase portrait of chaotic attractor when the follower offsets to the left $(O=20 \mathrm{~mm})$.
chaotic analysis in follower movement when the follower has detached from the cam. The system in Figure 8 has smooth orbit of follower displacement at (I.D. $=16 \mathrm{~mm}$ and 19 mm ) when the follower offsets to the left $(\mathrm{O}=20 \mathrm{~mm})$ at diverse cam speeds.


Figure 9.
Poincare' maps of chaotic attractor when the follower offsets to the left $(O=20 \mathrm{~mm})$.
The follower displacement is repeated itself based on the single black dots in phaseplane orbit. The chaotic analysis is detected based on the multi black dots in phaseplane orbit at (I.D. $=16 \mathrm{~mm}$ and 19 mm ) and diverse cam speeds as shown in Figure 9. SolidWorks software is used in the simulation.

## 6. Follower displacement

Figures 10 and 11 show the follower linear displacement against the time at different cam angular speeds when the camshaft offsets to the left ( $\mathrm{O}=40 \mathrm{~mm}$ ) and to the right $(\mathrm{O}=50 \mathrm{~mm})$ at $(\mathrm{I} . \mathrm{D} .=17 \mathrm{~mm})$ respectively. The follower stays in permanent contact when the cam starts spinning at ( $\mathrm{N}=200 \mathrm{rpm}$ ) and ( $\mathrm{N}=400 \mathrm{rpm}$ ), while the follower starts detaching from the cam at $(\mathrm{N}=1000 \mathrm{rpm})$ as shown in Figure 10. The follower also starts jumping a little bit higher from the cam at ( $\mathrm{N}=800 \mathrm{rpm}$ ) as shown in Figure 11.

MATLAB Code:
The code algorithm of phase-plane diagram and Poincare' map are added at the end of this chapter. The code is done using MATLAB software and as in below:

```
clear; clc; close all
SignalName = '100rpm.dat';
signal = load(SignalName);
signal = signal - min(signal);
% Poincare map
% original D=signal; % Read data
[x1max,t1max] = findpeaks(D(:,1));
```

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Figure 10.
Follower displacement against time when the follower offsets to the right $(O=50 \mathrm{~mm})$ at various cam speeds.


Figure 11.
Follower displacement against time when the follower offsets to the left $(O=40 \mathrm{~mm})$ at various cam speeds.

```
    Nmax = length(x1max); figure(1)
    subplot(1,2,1)
    for i=1:Nmax-1 plot(x1max(i),x1max(i+1),'ko','MarkerSize',5,'MarkerFa-
ceColor','k')
    hold on
    axis square
    xlabel('x_{max} '),ylabel('next x_{max}')
    grid on
    end
    %title('n = 100 rpm c = 1.5 mm') SignalName = '100rpmc2.dat';
    signal = load(SignalName);
    D = signal; % Read data
```

```
    [x1max,t1max] = findpeaks(D(:,1));
    Nmax = length(x1max); subplot(1,2,2)
    for i=1:Nmax-1 plot(x1max(i),x1max(i+1),'ko','MarkerSize',5,'MarkerFa-
ceColor','k')
    hold on
    axis square
    xlabel('x_{max} '),ylabel('next x_{max}')
    grid on
    end
    %title('n=100 rpm c = 2 mm') figure(2)
    aa = load('100rpm.dat');
    aa = aa - min(aa);
    plot(aa, gradient(aa));
```


## 7. Conclusions

In this article the chaotic motion of the follower response is considered in the presence of impact coefficient of restitution using SolidWorks program. The chaotic motion of the follower response is occurred due to the increase in cam speeds, follower's offsets, follower guides' clearances and impact coefficient of restitution. The value of Lyapunov exponent is increased with the increasing of embedding dimensions values. The positive local Lyapunov exponent depicts the transient state in nonlinear response of the follower in the presence of impact coefficient of restitution. Negative local Lyapunov exponent refers to the steady state in the follower motion. Some of the nonlinear response of the follower has periodic and chaotic motions at different time periods. The quantity of the black dots in Poincare' maps detects the chaotic analysis in follower movement when the follower has detached from the cam.

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## Conflict of interest

The author declares that he has no conflict of interest.

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## Chapter 6

# Decision Fusion for Large-Scale Sensor Networks with Nonideal Channels 

Yiwei Liao, Xiaojing Shen, Junfeng Wang and Yunmin Zhu


#### Abstract

Since there has been an increasing interest in the areas of Internet of Things (IoT) and artificial intelligence that often deals with a large number of sensors, this chapter investigates the decision fusion problem for large-scale sensor networks. Due to unavoidable transmission channel interference, we consider sensor networks with nonideal channels that are prone to errors. When the fusion rule is fixed, we present the necessary condition for the optimal sensor rules that minimize the Monte Carlo cost function. For the $K$-out-of- $L$ fusion rule chosen very often in practice, we analytically derive the optimal sensor rules. For general fusion rules, a Monte Carlo Gauss-Seidel optimization algorithm is developed to search for the optimal sensor rules. The complexity of the new algorithm is of the order of $O(L N)$ compared with $O\left(L N^{L}\right)$ of the previous algorithm that was based on Riemann sum approximation, where $L$ is the number of sensors and $N$ is the number of samples. Thus, the proposed method allows us to design the decision fusion rule for large-scale sensor networks. Moreover, the algorithm is generalized to simultaneously search for the optimal sensor rules and the optimal fusion rule. Finally, numerical examples show the effectiveness of the new algorithms for large-scale sensor networks with nonideal channels.


Keywords: decision fusion, multisensor detection, nonideal channels, Monte Carlo method, importance sampling

## 1. Introduction

Distributed detection has been an active research area in the past decades [1-7]. It involves the design of decision rules for the sensors ${ }^{1}$ and the fusion rule [8]. Early work on distributed detection mainly focused on conditionally independent sensor observations, such as [ $2,4,9,10$ ], and the resulting optimal sensor decision rules, as well as the fusion rule, were likelihood ratio tests (LRTs). Details on distributed detection with conditionally independent sensor observations can be seen in [1, 6, 7] and references therein.

[^0]In this chapter, we focus on conditionally dependent observations in sensor networks. In [5], the computational difficulty of obtaining the optimal sensor rules was shown by a rigorous mathematical approach. Some early progress was made on the derivation of sensor rules for the dependent observation case such as in [11-15]. More recently, a hierarchical conditional independence model was provided that was applicable to some specific classes of multisensor detection problems with dependent observations [16]. Copula-based distributed decision fusion methods have been proposed to deal with dependent observations in sensor networks, such as [17-19] and references therein. Given a fusion rule, Monte Carlo methods were proposed to reduce the computational complexity of deriving sensor decision rules with ideal channels in [20, 21], and the optimal sensor rules were obtained analytically for the $K$-out-of- $L$ fusion rule in [20].

Some works on the derivation of optimal fusion rules can be seen in [15, 22-24]. For some specific parallel network decision systems, a unified fusion rule was presented in [15]. Some further results on the problem are available in [25, 26]. In [27], the authors provided methods that search for the sensor rules and the fusion rule simultaneously by combining the methods of [2] and [15] in order to attain nearoptimal system performance.

The works discussed thus far assumed the availability of ideal channels in sensor networks. However, channel errors between the sensors and the fusion center are omnipresent in practical multisensor detection networks, and, therefore, studies on multisensor detection in the presence of nonideal channels have attracted some recent interest, such as in [8, 28-33]. Under the Neyman-Pearson criterion, the design of sensor rules in the presence of nonideal channels was addressed in [32]. The parallel fusion structure was extended by incorporating the fading channel layer and two alternative fusion schemes were presented based on fixed sensor rules in [28]. It was shown that the optimal sensor decision rule that minimizes the error probability at the fusion center is equivalent to a local LRT for independent sensor observations in [29]. Under NeymanPearson and Bayesian criteria, the work was generalized to dependent and noisy channels, respectively, in [8]. In [31], the authors considered the optimal sensor rules with channel errors via Riemann sum approximation under a given fusion rule for general dependent sensor observations. Although the method based on the Riemann sum approximation has been developed for dependent observations with channel errors, it is too computationally expensive to be of practical use in large-scale sensor networks.

In this chapter, a Monte Carlo importance sampling method is provided to reduce the computational complexity of multisensor detection fusion with channel errors. Based on the strong law of large numbers, the Bayesian cost function is approximated by its empirical average through the Monte Carlo importance sampling method. The main contributions of this chapter are listed below:

1. When the fusion rule is fixed, we derive a necessary condition for the optimal sensor rules that minimize the approximated Bayesian cost function. A Monte Carlo Gauss-Seidel optimization algorithm is developed and it is shown to be finitely convergent. The complexity of the new algorithm is shown to be of the order of $O(L N)$ compared with $O\left(L N^{L}\right)$ of the previous algorithm based on the Riemann sum approximation.
2. When the fusion rule is the $K$-out-of- $L$ rule, we prove that there exists an analytical form for the optimal sensor rules in the presence of nonideal channels. Thus, the proposed method allows us to design decision rules for large-scale sensor networks.
3. The Monte Carlo Gauss-Seidel optimization algorithm is extended to simultaneously search for the optimal sensor rules and the optimal fusion rule.

Numerical examples show the effectiveness of the new algorithms for large-scale sensor networks with dependent observations and channel errors.

The rest of this chapter is organized as follows: In Section 2, the parallel binary Bayesian detection network with channel errors is formulated and the Monte Carlo cost function is introduced. In Section 3, the necessary condition for the optimal sensor rules is presented. For the $K$-out-of- $L$ fusion rule, the analytical form for the optimal sensor rules is provided. In Section 4, the Monte Carlo Gauss-Seidel iterative algorithm and its convergence analysis are presented. The extension to search for the optimal sensor rules and the optimal fusion rule are simultaneously described in Section 5. Simulation results are provided in Section 6. Conclusions are contained in Section 7.

## 2. Preliminaries

### 2.1 Problem formulation

The $L$-sensor parallel Bayesian detection network structure with two hypotheses $H_{0}$ and $H_{1}$ in the presence of nonideal channels is considered (see Figure 1). Assume that $y_{1}, y_{2}, \ldots, y_{L}$ are sensor observations and the $j$ th sensor compresses the $n_{j}$ dimension vector observation $y_{j}$ to one bit: $I_{j}\left(y_{j}\right): \mathbb{R}^{n_{j}} \rightarrow\{0,1\}, j=1, \ldots, L$. For notational convenience, $n_{j}=1$ in the following description. The $L$ sensors transmit the compressed data to the fusion center and the fusion center makes the decision between $H_{0}$ and $H_{1}$. Since external interference and internal errors may occur, the channels are not reliable and the fusion center may not correctly receive the symbol $I_{j}$ sent by the $j$ th sensor. Let $I_{j}^{0}$ denote the received bit by the fusion center for $j=$ $1,2, \ldots, L$. Generally speaking, $I_{j}^{0}$ may not be equal to $I_{j}$. The definition and assumptions on channel errors (see e.g., $[29,31]$ ) are summarized below:

Definition 1: The channel errors between the $j$ th sensor and the fusion center are described as $P_{j}^{c e 1}=P\left(I_{j}^{0}=0 \mid I_{j}=1\right)$ and $P_{j}^{c e 0}=P\left(I_{j}^{0}=1 \mid I_{j}=0\right)$ for $j=1,2, \ldots, L$, where $P_{j}^{c e 1}$ is the probability of channel error when the $j$ th sensor sends 1 but the fusion center receives 0 , and $P_{j}^{c e 0}$ is the probability of channel error when the $j$ th sensor sends 0 but the fusion center receives 1 .


Figure 1.
The L-sensor parallel binary Bayesian detection network structure in the presence of nonideal channels.

Assumption 1: The probabilities of channel error are statistically independent of the hypotheses, namely $P\left(I_{j}^{0} \mid I_{j}, H_{\nu}\right)=P\left(I_{j}^{0} \mid I_{j}\right), \nu=0,1$.

Remark 1: Assumption 1 is due to the hierarchical structure based on the Markov property (see [29]).

Assumption 2: The channels that connect the sensors to the fusion center are independent, i.e., $P\left(I_{1}^{0}, I_{2}^{0}, \ldots, I_{L}^{0} \mid I_{1}, I_{2}, \ldots, I_{L}\right)=\prod_{j=1}^{L} P\left(I_{j}^{0} \mid I_{j}\right)$.

We consider the parallel binary Bayesian detection network with nonideal channels that is built on the above definition and assumptions. The final decision is made by the fusion center based on the received binary bits $\left(I_{1}^{0}, I_{2}^{0}, \ldots, I_{L}^{0}\right)$ from the $L$ sensors. From the definition of a general Bayesian cost function given in [25], the $L$-sensor binary Bayesian cost function with channel errors at the fusion center can be written as follows:

$$
\begin{gather*}
C\left(I_{1}^{0}\left(y_{1}\right), \ldots, I_{L}^{0}\left(y_{L}\right) ; F^{0}\right)=C_{00} P_{0} P\left(F^{0}=0 \mid H_{0}\right)+C_{01} P_{1} P\left(F^{0}=0 \mid H_{1}\right)  \tag{1}\\
+C_{10} P_{0} P\left(F^{0}=1 \mid H_{0}\right)+C_{11} P_{1} P\left(F^{0}=1 \mid H_{1}\right) \\
=c+a P\left(F^{0}=0 \mid H_{1}\right)-b P\left(F^{0}=0 \mid H_{0}\right) \tag{2}
\end{gather*}
$$

where $C_{\alpha \beta}, \alpha, \beta=0,1$ are suitable cost coefficients, $P_{0}$ and $P_{1}$ are the prior probabilities for the hypotheses $H_{0}$ and $H_{1}$, respectively, $F^{0}$ is the fusion rule, and $P\left(F^{0}=\mu \mid H_{\nu}\right), \mu, \nu=0,1$ denotes the conditional probability of the event that the fusion center decides in favor of hypothesis $\mu$ when the real hypothesis is $H_{\nu}$. The cost function (1) is simplified to (2) by defining $c=C_{10} P_{0}+C_{11} P_{1}, a=P_{1}\left(C_{01}-C_{11}\right)$, $b=P_{0}\left(C_{10}-C_{00}\right) . F^{0}$ is actually a function of the disjoint set of all possible binary messages $\left(I_{1}^{0}, I_{2}^{0}, \ldots, I_{L}^{0}\right)$. The received decisions are divided into two sets denoted as $H_{0}^{0}$ and $H_{1}^{0}$ which are given by

$$
\begin{aligned}
& H_{0}^{0}=\left\{\left(u_{1}^{0}, u_{2}^{0}, \ldots, u_{L}^{0}\right): F^{0}\left(\left(I_{1}^{0}, I_{2}^{0}, \ldots, I_{L}^{0}\right)\right)=0, I_{j}^{0}=u_{j}^{0}, u_{j}^{0}=0 / 1, j=1, \ldots, L\right\} \\
& H_{1}^{0}=\left\{\left(u_{1}^{0}, u_{2}^{0}, \ldots, u_{L}^{0}\right): F^{0}\left(\left(I_{1}^{0}, I_{2}^{0}, \ldots, I_{L}^{0}\right)\right)=1, I_{j}^{0}=u_{j}^{0}, u_{j}^{0}=0 / 1, j=1, \ldots, L\right\}
\end{aligned}
$$

Obviously, $H^{0}=\left\{\left(u_{1}^{0}, u_{2}^{0}, \ldots, u_{L}^{0}\right): I_{j}^{0}=u_{j}^{0}, u_{j}^{0}=0 / 1, j=1, \ldots, L\right\}=H_{0}^{0} \cup H_{1}^{0}$.
For any binary decisions $\left(I_{1}^{0}, I_{2}^{0}, \ldots, I_{L}^{0}\right)$ received by the fusion center, the original sensor decision bits before transmission are $\left(I_{1}, I_{2}, \ldots, I_{L}\right)$ and they consist of the set $H=\left\{\left(u_{1}, u_{2}, \ldots, u_{L}\right): I_{j}=u_{j}, u_{j}=0 / 1, j=1, \ldots, L\right\}$. Therefore, based on the law of total probability, the conditional probability formula, and Assumption 1:

$$
\begin{equation*}
P\left(F^{0}=0 \mid H_{\nu}\right)=\sum_{s^{0} \in H_{0}^{0}} P\left(D^{0} \mid H_{\nu}\right)=\sum_{s^{0} \in H_{0}^{0} s \in H} \sum_{i} P\left(D^{0} \mid D\right) P\left(D \mid H_{\nu}\right), \tag{3}
\end{equation*}
$$

where $D^{0}=\left(I_{1}^{0}, I_{2}^{0}, \ldots, I_{L}^{0}\right), s^{0}=\left(s^{0}(1), \ldots, s^{0}(L)\right), I_{j}^{0}=s^{0}(j)$, and $s^{0}(j)=0 / 1$ is a specific value of $I_{j}^{0}$; in the same way, $D=\left(I_{1}, I_{2}, \ldots, I_{L}\right), s=(s(1), \ldots, s(L)), I_{j}=s(j)$, and $s(j)=0 / 1$ is a specific value of $I_{j}$. Strictly speaking, we should use $P\left(D^{0}=s^{0} \mid H_{\nu}\right)$ to represent $P\left(D^{0} \mid H_{\nu}\right)$ and we use the latter for notational simplicity. It is similar to $P\left(D \mid H_{\nu}\right)$. Based on Assumption 2:

$$
\begin{equation*}
P\left(D^{0} \mid D\right)=\prod_{j=1}^{L} P\left(I_{j}^{0} \mid I_{j}\right) \tag{4}
\end{equation*}
$$

where for any $1 \leq j \leq L$

$$
\begin{equation*}
P\left(I_{j}^{0} \mid I_{j}\right)=\left(1-P_{j}^{c e 0}\right)\left(1-I_{j}^{0}\right)\left(1-I_{j}\right)+P_{j}^{c e 0} I_{j}^{0}\left(1-I_{j}\right)+\left(1-P_{j}^{c e 1}\right) I_{j}^{0} I_{j}+P_{j}^{c e 1}\left(1-I_{j}^{0}\right) I_{j} . \tag{5}
\end{equation*}
$$

Thus, the cost function (2) becomes

$$
\begin{gather*}
C\left(I_{1}^{0}\left(y_{1}\right), \ldots, I_{L}^{0}\left(y_{L}\right) ; F^{0}\right)=c+\sum_{s^{0} \in H_{0}^{0} \in H} P\left(D^{0} \mid D\right)\left[a P\left(D \mid H_{1}\right)-b P\left(D \mid H_{0}\right)\right]  \tag{6}\\
\triangleq C\left(I_{1}\left(y_{1}\right), \ldots, I_{L}\left(y_{L}\right) ; F^{0} ; P^{c e 0}, P^{c e 1}\right), \tag{7}
\end{gather*}
$$

where $P^{c e 0}=\left(P_{1}^{c e 0}, \ldots, P_{L}^{c e 0}\right), P^{c e 1}=\left(P_{1}^{c e 1}, \ldots, P_{L}^{c e 1}\right)$. Hence, the cost function now becomes a function of the sensor rules $\left(I_{1}, \ldots, I_{L}\right)$, the probabilities of channel errors $P^{c e 0}, P^{c e 1}$, and the fusion rule $F^{0}$. The goal of this chapter is to optimize the sensor rules and the fusion rule so as to minimize the cost function with known probabilities of channel errors.

We rewrite $a P\left(D \mid H_{1}\right)-b P\left(D \mid H_{0}\right)$ as follows:

$$
\begin{align*}
a P D\left|H_{1}-b P D\right| H_{0} & =\int_{\Omega_{s}} a p y_{1}, \ldots, y_{L}\left|H_{1}-b p y_{1}, \ldots, y_{L}\right| H_{0} d y_{1} \cdots d y_{L} \\
& =\int I_{\Omega_{s}}\left[a p y_{1}, \ldots, y_{L}\left|H_{1}-b p y_{1}, \ldots, y_{L}\right| H_{0}\right] d y_{1} \cdots d y_{L}, \tag{8}
\end{align*}
$$

where $\Omega_{s}=\left\{\left(y_{1}, \ldots, y_{L}\right): I_{1}\left(y_{1}\right)=s(1), \ldots, I_{L}\left(y_{L}\right)=s(L)\right\}, I_{\Omega_{s}}$ is an indicator function on $\Omega_{s}$, and the region of integration in (8) is the full space. Assume that $p\left(y_{1}, y_{2}, \ldots, y_{L} \mid H_{\nu}\right), \nu=0,1\left(\right.$ or $\left.p\left(y \mid H_{\nu}\right)\right)$ are the known conditional joint probability density functions. If not, we can learn the joint probability density functions from training data using copula functions (see, e.g., [17]). Note that $I_{1}\left(y_{1}\right), \ldots, I_{L}\left(y_{L}\right)$ are indicator functions and $s(j)=0 / 1, j=1, \ldots, L$,

$$
\begin{align*}
I_{\Omega_{s}} & =I_{\left\{\left(y_{1}, \ldots, y_{L}\right): I_{1}\left(y_{1}\right)=s(1), \ldots, I_{L}\left(y_{L}\right)=s(L)\right\}} \\
& =I_{\left\{y_{1}: I_{1}\left(y_{1}\right)=s(1)\right\}} \cdots I_{\left\{y_{L}: I_{L}\left(y_{L}\right)=s(L)\right\}}  \tag{9}\\
& =\left[\left(1-I_{1}\right)(1-s(1))+I_{1} s(1)\right] \cdots\left[\left(1-I_{L}\right)(1-s(L))+I_{L} s(L)\right] .
\end{align*}
$$

For simplicity, denote $Q_{j}\left(I_{j}\right)=\left(1-I_{j}\right)(1-s(j))+I_{j} s(j)$. Substituting (8) into (6),
$C\left(I_{1}\left(y_{1}\right), \ldots, I_{L}\left(y_{L}\right) ; F^{0} ; P^{c e 0}, P^{c e 1}\right)$
$=c+\sum_{s^{0} \in H_{0}^{0} \in H} \sum P\left(D^{0} \mid D\right) \cdot \int Q_{1}\left(I_{1}\right) \cdots Q_{L}\left(I_{L}\right)\left[a p\left(y \mid H_{1}\right)-b p\left(y \mid H_{0}\right)\right] d y=c+\int P_{H_{0}^{0}} \hat{L}(y) d y$,
where $P_{H_{0}^{0}}=\sum_{s^{0} \in H_{0}^{0}} \sum_{s \in H} P\left(D^{0} \mid D\right) Q_{1}\left(I_{1}\right) \cdots Q_{L}\left(I_{L}\right)$ and $\hat{L}(y)=a p\left(y \mid H_{1}\right)-$ $b p\left(y \mid H_{0}\right)$. Note that from the definition of $H_{0}^{0}, H_{1}^{0}, H^{0}$, and $H$, we have

$$
\begin{align*}
P_{H_{0}^{0}} & =\sum_{s^{0} \in H^{0}}\left[1-F^{0}\left(D^{0}\right)\right] \sum_{s \in H} P\left(D^{0} \mid D\right) Q_{1}\left(I_{1}\right) \cdots Q_{L}\left(I_{L}\right) \\
& =\sum_{k^{\prime}=1}^{2^{L}} \sum_{k=1}^{2^{L}}\left[1-F^{0}\left(s_{k^{\prime}}\right)\right] P\left(s_{k^{\prime}} \mid s_{k}\right) \cdot \prod_{j=1}^{L}\left\{\left[1-I_{j}\left(y_{j}\right)\right]\left[1-s_{k}(j)\right]+I_{j}\left(y_{j}\right) s_{k}(j)\right\}, \tag{11}
\end{align*}
$$

where $s_{k^{\prime}}$ is the element of $H^{0}$ and $s_{k}$ is the element of $H . F^{0}\left(D^{0}\right)=0 / 1$ is used in the first equality. The second equality holds since there are $2^{L}$ elements in both $H$ and $H^{0}$.

### 2.2 Monte Carlo cost function

An essential difficulty of the Bayesian cost function (10) is the required high dimensional integration when dealing with large-scale sensor networks. Monte Carlo importance sampling is an attractive method to deal with this problem. In this subsection, we approximate the Bayesian cost function (10) by the Monte Carlo importance sampling method (see, e.g., [34, 35]). According to (10),

$$
\begin{align*}
& C\left(I_{1}\left(y_{1}\right), \ldots, I_{L}\left(y_{L}\right) ; F^{0} ; P^{c e 0}, P^{c c 1}\right) \\
& \begin{array}{c}
=c+\int P_{H_{0}^{0}}\left(I_{1}\left(y_{1}\right), \ldots, I_{L}\left(y_{L}\right) ; F^{0} ; P^{c e 0}, P^{c c 1}\right) \frac{\hat{L}(y)}{g(y)} \cdot g(y) d y \\
=\mathbb{E}_{g} \frac{P_{H_{0}^{0}}(Y) \hat{L}(Y)}{g(Y)}+c
\end{array} \tag{12}
\end{align*}
$$

where $y=\left(y_{1}, y_{2}, \ldots, y_{L}\right)$, and $g(y)$ is a given importance sampling density such that (12) is well-defined (i.e., $g(y)>0$ ). In (13), the expectation is taken with respect to the importance sampling density $g$. Consequently, assume that $N$ samples $Y_{1}, \ldots, Y_{N}$ are generated from the density $g$, that is, $Y \sim g(y)$, where $Y_{i}=$ $\left[Y_{i 1}, Y_{i 2}, \ldots, Y_{i L}\right]$. Then

$$
\begin{gather*}
C\left(I_{1}\left(y_{1}\right), \ldots, I_{L}\left(y_{L}\right) ; F^{0} ; P^{c e 0}, P^{c e 1}\right) \approx \frac{1}{N} \sum_{i=1}^{N} \frac{P_{H_{0}^{0}}\left(Y_{i 1}, Y_{i 2}, \ldots, Y_{i L}\right) \hat{L}\left(Y_{i}\right)}{g\left(Y_{i}\right)}+c  \tag{14}\\
\triangleq C_{M C}\left(I_{1}\left(y_{1}\right), \ldots, I_{L}\left(y_{L}\right) ; F^{0} ; P^{c e 0}, P^{c e 1}, N\right) . \tag{15}
\end{gather*}
$$

Based on the strong law of large numbers, the expectation (13) can be approximated by the empirical average (14). Denote (14), namely the Monte Carlo cost function, as $C_{M C}\left(I_{1}\left(y_{1}\right), \ldots, I_{L}\left(y_{L}\right) ; F^{0} ; P^{c e 0}, P^{c e 1}, N\right)$. The optimal importance sampling density is $g\left(y_{1}, y_{2}, \ldots, y_{L}\right) \propto\left|P_{H_{0}^{0}} \hat{L}\left(y_{1}, y_{2}, \ldots, y_{L}\right)\right|$ (see, e.g., $[34,35]$ ).

The initial goal is to minimize the Bayesian cost function (10). Instead, we can minimize the Monte Carlo cost function (15) by selecting a set of optimal sensor rules $I_{1}\left(y_{1}\right), I_{2}\left(y_{2}\right), \ldots, I_{L}\left(y_{L}\right)$ and an optimal fusion rule $F^{0}$. In this manner, the highdimensional integration problem is converted to a problem where we need to deal with the single summation objective function for large-scale sensor networks. Thus, for dependent observations with channel errors, the computational complexity is reduced significantly by the Monte Carlo importance sampling method. In the following sections, we assume that the samples drawn from the importance sampling
density are fixed so that $C_{M C}\left(I_{1}, \ldots, I_{L} ; F^{0} ; P^{c e 0}, P^{c e 1}, N\right)$ does not have any randomness, since only deterministic decision rules are considered in this chapter.

## 3. A necessary condition for the optimal sensor rules

In this section, when the fusion rule is fixed, we derive a necessary condition for the optimal sensor rules that minimize the Monte Carlo cost function. First, we need some equivalent transformations for $P_{H_{0}^{0}}$. Then based on the transformations, the necessary condition can be obtained. At the same time, an analytical result is obtained when the fusion rule is the $K$-out-of- $L$ rule.

### 3.1 Necessary condition

First, we need some equivalent transformations for $P_{H_{0}^{0}}$.
Lemma $1 P_{H_{0}^{0}}$ can be rewritten as follows:

$$
\begin{align*}
P_{H_{0}^{0}} & \triangleq\left[1-I_{j}\left(y_{j}\right)\right] P_{j 1}\left(I_{1}\left(y_{1}\right), \ldots, I_{j-1}\left(y_{j-1}\right), I_{j+1}\left(y_{j+1}\right), \ldots, I_{L}\left(y_{L}\right) ; F^{0} ; P^{c e 0}, P^{c e 1}\right) \\
& +P_{j 2}\left(I_{1}\left(y_{1}\right), \ldots, I_{j-1}\left(y_{j-1}\right), I_{j+1}\left(y_{j+1}\right), \ldots, I_{L}\left(y_{L}\right) ; F^{0} ; P^{c e 0}, P^{c e 1}\right) \tag{16}
\end{align*}
$$

where for $j=1,2, \ldots, L$,

$$
\begin{array}{r}
P_{j 1}(\cdot) \triangleq \sum_{k^{\prime}=1}^{2^{L}}\left[1-F\left(s_{k^{\prime}}\right)\right]\left(1-P_{j}^{c e 0}-P_{j}^{c e 1}\right)\left(1-2 s_{k^{\prime}}(j)\right) P_{m \neq j}, \\
P_{j 2}(\cdot) \triangleq \sum_{k^{\prime}=1}^{2^{L}}\left[1-F\left(s_{k^{\prime}}\right)\right]\left[s_{k^{\prime}}(j)+P_{j}^{c e 1}\left(1-2 s_{k^{\prime}}(j)\right)\right] P_{m \neq j}, \\
P_{m \neq j} \triangleq \prod_{m, m \neq j}^{L}\left\{\left(1-P_{m}^{c e 0}\right)\left(1-s_{k^{\prime}}(m)\right)\left(1-I_{m}\right)+P_{m}^{c e 0} s_{k^{\prime}}(m)\left(1-I_{m}\right)\right.  \tag{19}\\
\left.+\left(1-P_{m}^{c e 1}\right) s_{k^{\prime}}(m) I_{m}+P_{m}^{c e 1}\left(1-s_{k^{\prime}}(m)\right) I_{m}\right\} .
\end{array}
$$

Proof: If $s_{k}(m)=I_{m}\left(y_{m}\right)$ for all $m=1, \ldots, L$, then the continued product term $\prod_{m=1}^{L}\left\{\left[1-I_{m}\left(y_{m}\right)\right]\left[1-s_{k}(m)\right]+I_{m}\left(y_{m}\right) s_{k}(m)\right\}=1$ in $P_{H_{0}^{0}}$. Otherwise, it is 0 . Thus, $P_{H_{0}^{0}}$ can be rewritten as $P_{H_{0}^{0}}=\sum_{k^{\prime}=1}^{2^{L}}\left[1-F\left(s_{k^{\prime}}\right)\right] P\left(s_{k^{\prime}} \mid\left(I_{1}, I_{2}, \ldots, I_{L}\right)\right)$, where the terms that equal zero are omitted and $\left.P\left(s_{k^{\prime}} \mid\left(I_{1}, I_{2}, \ldots, I_{L}\right)\right)=\prod_{j=1}^{L} P\left(s_{k^{\prime}} j\right) \mid I_{j}\right)$. Recalling the conditional probability formula (5), we rewrite $P\left(s_{k^{\prime}}(j) \mid I_{j}\right)$ as $\left.P\left(s_{k^{\prime}} j\right) \mid I_{j}\right)=$ $\left[1-I_{j}\right]\left(1-P_{j}^{c e 0}-P_{j}^{c e 1}\right)\left(1-2 s_{k^{\prime}}(j)\right)+s_{k^{\prime}}(j)+P_{j}^{c e 1}\left(1-2 s_{k^{\prime}}(j)\right)$. Based on these transformations, $P_{H_{0}^{0}}$ can be decomposed as (16).

Remark 2: Note that $P_{j 1}(\cdot)$ and $P_{j 2}(\cdot)$ are both independent of $I_{j}\left(y_{j}\right)$ for $j=1, \ldots, L$. In addition, they can also be applied in the Riemann sum approximation (see, e.g., [31]).

Compared with [36], the sum of $2^{L}$ terms about $s_{k}$ is eliminated and it greatly reduces the computational time. In addition, the expression for $P_{j 1}(\cdot)$ given in (17) is also a key equation in the following results:

Substituting the transformations (16) into (15), we obtain

$$
\begin{align*}
& C_{M C}\left(I_{1}\left(y_{1}\right), \ldots, I_{L}\left(y_{L}\right) ; F^{0} ; P^{c c 0}, P^{c c 1}, N\right) \\
& =c+\frac{1}{N} \sum_{i=1}^{N}\left\{\left[1-I_{j}\left(Y_{i j}\right)\right] P_{j 1}\left(I_{1}\left(Y_{i 1}\right), \ldots, I_{j-1}\left(Y_{i(j-1)}\right), I_{j+1}\left(Y_{i(j+1)}\right), \ldots, I_{L}\left(Y_{i L}\right) ; F^{0} ; P^{c e 0}, P^{c e 1}\right)\right. \\
& \left.+P_{j 2}\left(I_{1}\left(Y_{i 1}\right), \ldots, I_{j-1}\left(Y_{i j-1)}\right), I_{j+1}\left(Y_{i j+1)}\right), \ldots, I_{L}\left(Y_{i L}\right) ; F^{0} ; P^{c e 0}, P^{c e 1}\right)\right\} \cdot \frac{\hat{L}\left(Y_{i}\right)}{g\left(Y_{i}\right)}, \tag{20}
\end{align*}
$$

where $Y_{i}=\left(Y_{i 1}, Y_{i 2}, \ldots, Y_{i L}\right)$. According to (20), the necessary condition for the optimal sensor rules that minimize $C_{M C}\left(I_{1}\left(y_{1}\right), \ldots, I_{L}\left(y_{L}\right) ; F^{0} ; P^{c e 0}, P^{c c 1}, N\right)$ is stated in the following lemma:

Lemma 2: Let $\left\{I_{1}\left(y_{1}\right), \ldots, I_{L}\left(y_{L}\right)\right\}$ be the set of optimal sensor rules, i.e., they minimize $C_{M C}\left(I_{1}\left(y_{1}\right), \ldots, I_{L}\left(y_{L}\right) ; F^{0} ; P^{c e 0}, P^{c e 1}, N\right)$ in the parallel Bayesian detection network, then they must satisfy the following equations:

$$
\begin{gather*}
I_{1}\left(Y_{i 1}\right)=I\left[P_{11}\left(I_{2}\left(Y_{i 2}\right), I_{3}\left(Y_{i 3}\right), \ldots, I_{L}\left(Y_{i L}\right) ; F^{0} ; P^{c e 0}, P^{c e 1}\right) \cdot \hat{L}\left(Y_{i}\right)\right],  \tag{21}\\
I_{2}\left(Y_{i 2}\right)=I\left[P_{21}\left(I_{1}\left(Y_{i 1}\right), I_{3}\left(Y_{i 3}\right), \ldots, I_{L}\left(Y_{i L}\right) ; F^{0} ; P^{c e 0}, P^{c e 1}\right) \cdot \hat{L}\left(Y_{i}\right)\right],  \tag{22}\\
\ldots \ldots  \tag{23}\\
I_{L}\left(Y_{i L}\right)=I\left[P_{L 1}\left(I_{1}\left(Y_{i 1}\right), I_{2}\left(Y_{i 2}\right), \ldots, I_{L-1}\left(Y_{i(L-1)}\right) ; F^{0} ; P^{c e 0}, P^{c e 1}\right) \cdot \hat{L}\left(Y_{i}\right)\right],
\end{gather*}
$$

where $P_{j 1}(\cdot)$ are defined by (17) and $I[\cdot]$ is an indicator function defined as follows:

$$
I[x]= \begin{cases}1, & \text { if } x \geq 0  \tag{24}\\ 0, & \text { if } \quad x<0\end{cases}
$$

Proof: Note that both $P_{j 1}(\cdot)$ and $P_{j 2}(\cdot)$ are independent of $I_{j}\left(y_{j}\right)$ for $j=1, \ldots, L$. If $I_{1}\left(Y_{i 1}\right)$ minimizes the Monte Carlo cost function under the given $I_{2}\left(Y_{i 2}\right), \ldots, I_{L}\left(Y_{i L}\right)$, we only need to minimize the first term of the summation in (20), that is, $\left[1-I_{1}\left(Y_{i 1}\right)\right] P_{11}\left(I_{2}\left(Y_{i 2}\right), I_{3}\left(Y_{i 3}\right), \ldots, I_{L}\left(Y_{i L}\right) ; F^{0} ; P^{c e 0}, P^{c e 1}\right) \frac{\hat{L}\left(Y_{i}\right)}{g\left(Y_{i}\right)}$. Note that the value of $I_{1}\left(Y_{i 1}\right)$ is 0 or 1 and $g(y)$ is well defined, that is, $g\left(Y_{i}\right)>0, I_{1}\left(Y_{i 1}\right)$ should be equal to 1 when $P_{11}\left(I_{2}\left(Y_{i 2}\right), I_{3}\left(Y_{i 3}\right), \ldots, I_{L}\left(Y_{i L}\right) ; F^{0} ; P^{c c 0}, P^{c e 1}\right) \hat{L}\left(Y_{i}\right) \geq 0$ for $i=1, \ldots, N$, otherwise it should be equal to 0 . Therefore, we obtain (21) by the definition of $I[x]$ in (24). Similarly, we obtain (22) and (23) by minimizing (20).

### 3.2 An analytical result for the $K$-out-of- $L$ rule

When the fusion rule is a $K$-out-of- $L$ rule, we would obtain an analytical result in the presence of nonideal channels. It is described as follows:

Theorem 1.1: If the fusion rule is a $K$-out-of- $L$ rule and the probabilities of channel errors are less than 0.5 (i.e., $0<P_{j}^{c e 0}<0.5,0<P_{j}^{c e 1}<0.5$ ) for each channel, the optimal sensor rules are $I_{j}\left(Y_{i j}\right)=I\left[\hat{L}\left(Y_{i}\right)\right]$ for $i=1, \ldots, N$ and $j=1, \ldots, L$.

Proof: From Lemma 1, we know

$$
\begin{align*}
P_{j 1}(\cdot) & =\sum_{k^{\prime}=1}^{2^{L}}\left[1-F\left(s_{k^{\prime}}\right)\right] P_{m \neq j} \cdot\left(1-P_{j}^{c e 0}-P_{j}^{c e 1}\right)\left[1-2 s_{k^{\prime}}(j)\right] \\
& =\left(1-P_{j}^{c e 0}-P_{j}^{c c 1}\right) \sum_{k^{\prime}=1}^{2^{L-1}}\left[\left(1-F\left(s_{k^{\prime}} \mid s_{k^{\prime}}(j)=0\right)\right)-\left(1-F\left(s_{k^{\prime}} \mid s_{k^{\prime}}(j)=1\right)\right)\right] P_{m \neq j} \\
& =\left(1-P_{j}^{c e 0}-P_{j}^{c e 1}\right) \sum_{k^{\prime}=1}^{2^{L-1}}\left[F\left(s_{k^{\prime}} \mid s_{k^{\prime}}(j)=1\right)-F\left(s_{k^{\prime}} \mid s_{k^{\prime}}(j)=0\right)\right] P_{m \neq j} \tag{25}
\end{align*}
$$

Since $0<P_{j}^{c e 0}<0.5,0<P_{j}^{c e 1}<0.5$, we have $1-P_{j}^{c e 0}-P_{j}^{c e 1}>0$. Obviously, $P_{m \neq j}>0$ holds from its definition. If $\left[F\left(s_{k^{\prime}} \mid s_{k^{\prime}}(j)=1\right)-F\left(s_{k^{\prime}} \mid s_{k^{\prime}}(j)=0\right)\right] \geq 0, P_{j 1}(\cdot) \geq 0$ can be derived. When the fusion rule is a $K$-out-of- $L$ rule, $F\left(s_{k^{\prime}}\right)=I\left[\sum_{j=1}^{L} s_{k^{\prime}}(j)-K\right]$. Thus,

$$
\begin{aligned}
& F\left(s_{k^{\prime}} \mid s_{k^{\prime}}(j)=1\right)=I\left[\sum_{m=1, m \neq j}^{L} s_{k^{\prime}}(m)+1-K\right], \\
& F\left(s_{k^{\prime}} \mid s_{k^{\prime}}(j)=0\right)=I\left[\sum_{m=1, m \neq j}^{L} s_{k^{\prime}}(m)+0-K\right] .
\end{aligned}
$$

If $\sum_{m=1, m \neq j}^{L} s_{k^{\prime}}(m)+0-K \geq 0$, then $\sum_{m=1, m \neq j}^{L} s_{k^{\prime}}(m)+1-K \geq 0$, and we can get that $F\left(s_{k^{\prime}} s_{k^{\prime}}(j)=1\right)-F\left(s_{k^{\prime}} \mid s_{k^{\prime}}(j)=0\right)=0$. If $\sum_{m=1, m \neq j}^{L} s_{k^{\prime}}(m)+0-K<0$, then $F\left(s_{k^{\prime}} \mid s_{k^{\prime}}(j)=1\right)-F\left(s_{k^{\prime}} \mid s_{k^{\prime}}(j)=0\right) \geq 0$. In a word, $F\left(s_{k^{\prime}} \mid s_{k^{\prime}}(j)=1\right)-F\left(s_{k^{\prime}} \mid s_{k^{\prime}}(j)=0\right) \geq 0$ is derived, thus $P_{j 1} \geq 0$. It is easy to find a $s_{k^{\prime}}(m), m \neq j$ so that $\sum_{m=1, m \neq j}^{L} s_{k^{\prime}}(m)+0=$ $K-1$ and $\sum_{m=1, m \neq j}^{L} s_{k^{\prime}}(m)+1=K$. Thus, there must exist $F\left(s_{k^{\prime}} \mid s_{k^{\prime}}(j)=1\right)-$ $F\left(s_{k^{\prime}} \mid s_{k^{\prime}}(j)=0\right)>0$. Therefore, the $P_{j 1}>0$ is derived. Recalling the necessary condition for the optimal sensor rules, that is, $I\left(Y_{i j}\right)=I\left[P_{j 1}(\cdot) \cdot \hat{L}\left(Y_{i}\right)\right]$, the proof is completed.

Remark 3: The $K$-out-of- $L$ rule counts the number of sensors that vote in favor of $H_{1}$ and compares it with a given threshold $K$ [37]. It is also referred to as the counting rule or voting rule and is widely used in the practical decision fusion area [38, 39]. It encompasses a general class of fusion rules such as AND, OR, and Majority Boolean fusion rules [40]. The reason we assume that the probabilities of channel errors are less than 0.5 is based on practical considerations. If the probabilities of channel errors are greater than or equal to 0.5 , the channel is totally unreliable and the performance is not better than a random decision. Obviously, the analytical solution is very efficient to tackle large-scale sensor networks with dependent observations and channel errors.

## 4. Monte Carlo Gauss-Seidel iterative algorithm and its convergence

For general fusion rules that do not have the form of a $K$-out-of- $L$ rule, an efficient algorithm can be obtained that is inspired by Lemma 2. Next, we present a Monte Carlo Gauss-Seidel iterative algorithm and derive its convergence, when the fusion rule is fixed.

### 4.1 Monte Carlo Gauss-Seidel iterative algorithm

Based on the fixed-point type necessary condition given in Lemma 2, the Monte Carlo Gauss-Seidel iterative algorithm is presented in Algorithm 1.

Algorithm 1: Optimization of the sensor rules.
Given the fusion rule $F^{0}$ :

- Step 1: Generate $N$ samples: $Y_{1}, \ldots, Y_{N} \sim g(y)$, where $g(y)$ is an importance sampling density and $Y_{i}=\left[Y_{i 1}, Y_{i 2}, \ldots, Y_{i L}\right]$.
- Step 2: Initialize the $L$ sensor rules, for $j=1,2, \ldots, L$ and $i=1, \ldots, N$,

$$
\begin{equation*}
I_{j}^{(0)}\left(Y_{i j}\right)=0 / 1 \tag{26}
\end{equation*}
$$

- Step 3: Iteratively search for the $L$ sensor rules until a termination criterion in Step 4 is satisfied. The $n+1$ th iteration is given as follows: for $i=1, \ldots, N$,

$$
\begin{gather*}
I_{1}^{(n+1)}\left(Y_{i 1}\right)=I\left[P_{11}\left(I_{2}^{(n)}\left(Y_{i 2}\right), I_{3}^{(n)}\left(Y_{i 3}\right), \ldots, I_{L}^{(n)}\left(Y_{i L}\right) ; F^{0} ; P^{c e 0}, P^{c e 1}\right) \hat{L}\left(Y_{i}\right],\right.  \tag{27}\\
I_{2}^{(n+1)}\left(Y_{i 2}\right)=I\left[P_{21}\left(I_{1}^{(n+1)}\left(Y_{i 1}\right), I_{3}^{(n)}\left(Y_{i 3}\right), \ldots, I_{L}^{(n)}\left(Y_{i L}\right) ; F^{0} ; P^{c e 0}, P^{c e 1}\right) \hat{L}\left(Y_{i}\right)\right],  \tag{28}\\
I_{L}^{(n+1)}\left(Y_{i L}\right)=I\left[P_{L 1}\left(I_{1}^{(n+1)}\left(Y_{i 1}\right), \ldots, I_{L-1}^{(n+1)}\left(Y_{i(L-1)}\right) ; F^{0} ; P^{c e 0}, P^{c e 1}\right) \hat{L}\left(Y_{i}\right)\right], \tag{29}
\end{gather*}
$$

- Step 4: For $i=1, \ldots, N$, the termination criterion of the iteration process is

$$
\begin{gather*}
I_{1}^{(n+1)}\left(Y_{i 1}\right)=I_{1}^{(n)}\left(Y_{i 1}\right), \\
I_{2}^{(n+1)}\left(Y_{i 2}\right)=I_{2}^{(n)}\left(Y_{i 2}\right),  \tag{30}\\
\ldots \ldots \\
I_{L}^{(n+1)}\left(Y_{i L}\right)=I_{L}^{(n)}\left(Y_{i L}\right) .
\end{gather*}
$$

Remark 4: When we obtain $I_{1}\left(Y_{i 1}\right)$ for $i=1, \ldots, N$, we can compress $y_{1}$ by defining $I_{1}\left(y_{1}\right)=I_{1}\left(Y_{i 1}\right)$ if the distance $\left\|y_{1}-Y_{i 1}\right\| \leq\left\|y_{1}-Y_{i^{\prime} 1}\right\|$ for all $i^{\prime} \neq i$. In the same way, we can compress $y_{j}$ for $j=2, \ldots, L$. In fact, the method is to find one nearest neighbor of $y_{j}$ for $j=1, \ldots, L$ and use the corresponding compression rule. Moreover, we can utilize the k-nearest neighbor (knn) to compress $y_{j}$ (see more in [41]).

Remark 5: The main computation burden of Algorithm 1 is included in (27)-(29). If we let the number of discretized points $N_{1}=N_{2}=\ldots=N_{L}=N$ in [31], then $P_{j 1}(\cdot) \hat{L}\left(Y_{i}\right), j=1, \ldots, L$, and $i=1, \ldots, N$ are computed $L N^{L}$ times for each iteration, as in [31]. But they only need to be computed $L N$ times in Algorithm 1. Thus, the computational complexity of Algorithm 1, i.e., $O(L N)$ is much less than that in [31], that is, $O\left(L N^{L}\right)$. It is more efficient to tackle large-scale sensor networks with dependent observations and channel errors.

### 4.2 Convergence of the iterative algorithm

Now, we show that Algorithm 1 must converge to a stationary point and the algorithm cannot oscillate infinitely, that is, it terminates after a finite number of iterations.

Lemma 3: Given the fusion rule $F^{0}$, for any initial values $\left(I_{1}^{(0)}, \ldots, I_{L}^{(0)}\right)$ in (26), $C_{M C}\left(I_{1}^{(n)}, \ldots, I_{j}^{(n)}, I_{j+1}^{(n)}, \ldots, I_{L}^{(n)} ; F^{0} ; P^{c e 0}, P^{c e 1}, N\right)$ must converge to a stationary point after a finite number of iterations.

Proof: For $j=1, \ldots, L$, we denote $C_{M C}(20)$ in the $n+1$ th iteration process by

$$
\begin{align*}
& C_{M C}\left(I_{1}^{(n+1)}, \ldots, I_{j}^{(n+1)}, I_{j+1}^{(n)}, \ldots, I_{L}^{(n)} ; F^{0} ; P^{c e 0}, P^{c e 1}, N\right) \\
& =\frac{1}{N} \sum_{i=1}^{N}\left\{[ 1 - I _ { j } ^ { ( n + 1 ) } ( Y _ { i j } ) ] P _ { j 1 } \left(I_{1}^{(n+1)}\left(Y_{i 1}\right), \ldots, I_{j-1}^{(n+1)}\left(Y_{i(j-1)}\right), I_{j+1}^{(n)}\left(Y_{i(j+1)}\right), \ldots, I_{L}^{(n)}\left(Y_{i L}\right) ; F^{0} ;\right.\right. \\
& \left.\left.\quad P^{c e 0}, P^{c e 1}, N\right)+P_{j 2}\left(I_{1}^{(n+1)}\left(Y_{i 1}\right), \ldots, I_{j-1}^{(n+1)}\left(Y_{i(j-1)}\right), I_{j+1}^{(n)}\left(Y_{i(j+1)}\right), \ldots, I_{L}^{(n)}\left(Y_{i L}\right) ; F^{0} ; P^{c e 0}, P^{c e 1}, N\right)\right\} \\
& \quad \cdot \frac{\hat{L}\left(Y_{i}\right)}{g\left(Y_{i}\right)}+c . \tag{31}
\end{align*}
$$

Similarly, we denote the $(n+1)$ th iteration process of the iterative items $P_{j 1}(\cdot) \hat{L}(\cdot)$ in (27)-(29) by

$$
\begin{equation*}
G_{j}^{i}=P_{j 1}\left(I_{1}^{(n+1)}\left(Y_{i 1}\right), \ldots, I_{j-1}^{(n+1)}\left(Y_{i(j-1)}\right), I_{j+1}^{(n)}\left(Y_{i(j+1)}\right), \ldots, I_{L}^{(n)}\left(Y_{i L}\right) ; F^{0} ; P^{c c 0}, P^{c e 1}, N\right) \hat{L}\left(Y_{i}\right), \tag{32}
\end{equation*}
$$

for $i=1, \ldots, N$ and $j=1, \ldots, L$. Plugging $G_{j}^{i}$ into (31), we know

$$
\begin{equation*}
C_{M C}\left(I_{1}^{(n+1)}, \ldots, I_{j}^{(n+1)}, I_{j+1}^{(n)}, \ldots, I_{L}^{(n)} ; F^{0} ; P^{c e 0}, P^{c e 1}, N\right)=\frac{1}{N} \sum_{i=1}^{N} \frac{\left[1-I_{j}^{(n+1)}\left(Y_{i j}\right)\right]}{g\left(Y_{i}\right)} G_{j}^{i}+C_{j}^{i}, \tag{33}
\end{equation*}
$$

where $C_{j}^{i}=c+\frac{1}{N} \sum_{i=1}^{N} P_{j 2}\left(I_{1}^{(n+1)}\left(Y_{i 1}\right), \ldots, I_{j-1}^{(n+1)}\left(Y_{i(j-1)}\right), I_{j+1}^{(n)}\left(Y_{i(j+1)}\right), \ldots, I_{L}^{(n)}\left(Y_{i L}\right) ; F^{0}\right.$; $\left.P^{c e 0}, P^{c e 1}, N\right) \frac{\hat{L}\left(Y_{i}\right)}{g\left(Y_{i}\right)}$ is independent of $I_{j}^{(n)}$ and $I_{j}^{(n+1)}$. Splitting $1-I_{j}^{(n+1)}\left(Y_{i j}\right)$ into two terms, we obtain

$$
\begin{align*}
& C_{M C}\left(I_{1}^{(n+1)}, \ldots, I_{j}^{(n+1)}, I_{j+1}^{(n)}, \ldots, I_{L}^{(n)} ; F^{0} ; P^{c e 0}, P^{c e 1}, N\right) \\
& \quad=\frac{1}{N} \sum_{i=1}^{N} \frac{\left[1-I_{j}^{(n)}\left(Y_{i j}\right)\right]+\left[I_{j}^{(n)}\left(Y_{i j}\right)-I_{j}^{(n+1)}\left(Y_{i j}\right)\right]}{g\left(Y_{i}\right)} G_{j}^{i}+C_{j}^{i} \\
& \quad=\frac{1}{N} \sum_{i=1}^{N} \frac{\left[1-I_{j}^{(n)}\left(Y_{i j}\right)\right]}{g\left(Y_{i}\right)} G_{j}^{i}+C_{j}^{i}+\frac{1}{N} \sum_{i=1}^{N} \frac{\left[I_{j}^{(n)}\left(Y_{i j}\right)-I_{j}^{(n+1)}\left(Y_{i j}\right)\right]}{g\left(Y_{i}\right)} G_{j}^{i}  \tag{34}\\
& \quad=C_{M C}\left(I_{1}^{(n+1)}, \ldots, I_{j-1}^{(n+1)}, I_{j}^{(n)}, \ldots, I_{L}^{(n)} ; F^{0} ; P^{c e 0}, P^{c e 1}, N\right)+D_{j}^{(n+1)}
\end{align*}
$$

where

$$
\begin{equation*}
D_{j}^{(n+1)}=\frac{1}{N} \sum_{i=1}^{N} \frac{\left[I_{j}^{(n)}\left(Y_{i j}\right)-I_{j}^{(n+1)}\left(Y_{i j}\right)\right]}{g\left(Y_{i}\right)} G_{j}^{i} \tag{35}
\end{equation*}
$$

Note that (27)-(29) imply that $I_{j}^{(n+1)}\left(Y_{i j}\right)=0$ if and only if $G_{j}^{i}<0$ and $I_{j}^{(n+1)}\left(Y_{i j}\right)=1$ if and only if $G_{j}^{i} \geq 0$ for $i=1, \ldots, N, j=1, \ldots, L$. It means

$$
\begin{equation*}
\left[I_{j}^{(n)}\left(Y_{i j}\right)-I_{j}^{(n+1)}\left(Y_{i j}\right)\right] G_{j}^{i} \leq 0 \tag{36}
\end{equation*}
$$

Thus, for $\forall i, j$

$$
\begin{equation*}
\left[I_{j}^{(n)}\left(Y_{i j}\right)-I_{j}^{(n+1)}\left(Y_{i j}\right)\right] G_{j}^{i} / g\left(Y_{i}\right) \leq 0 \tag{37}
\end{equation*}
$$

where the inequality (35) holds since $g(\cdot)$ is well-defined (i.e., $g(\cdot)>0$ ). Substituting (35) into (33) yields $D_{j}^{(n+1)} \leq 0$. Thus, for $\forall j \leq L$,

$$
\begin{align*}
& C_{M C}\left(I_{1}^{(n+1)}, \ldots, I_{j}^{(n+1)}, I_{j+1}^{(n)}, \ldots, I_{L}^{(n)} ; F^{0} ; P^{c e 0}, P^{c c 1}, N\right) \\
& \leq C_{M C}\left(I_{1}^{(n+1)}, \ldots, I_{j-1}^{(n+1)}, I_{j}^{(n)}, \ldots, I_{L}^{(n)} ; F^{0} ; P^{c e 0}, P^{c e 1}, N\right) \tag{38}
\end{align*}
$$

Furthermore,

$$
\begin{align*}
& C_{M C}\left(I_{1}^{(n+1)}, I_{2}^{(n+1)}, \ldots, I_{L}^{(n+1)} ; F^{0} ; P^{c e 0}, P^{c e 1}, N\right) \\
& \leq C_{M C}\left(I_{1}^{(n)}, I_{2}^{(n)}, \ldots, I_{L}^{(n)} ; F^{0} ; P^{c e 0}, P^{c e 1}, N\right) \tag{39}
\end{align*}
$$

It means $C_{M C}$ is nonincreasing. Note that $C_{M C}\left(I_{1}^{(n)}, I_{2}^{(n)}, \ldots, I_{L}^{(n)} ; F^{0} ; P^{c e 0}, P^{c e 1}, N\right)$ is a finite value. We conclude that it must converge to a stationary point after a finite number of iterations.

Theorem 1.2: Given the fusion rule $F^{0}$, the sensor rules $I_{1}^{(n)}, I_{2}^{(n)}, \ldots, I_{L}^{(n)}$ are finitely convergent, i.e., Algorithm 1 converges after a finite number of iterations.

Proof: By Lemma 3, $C_{M C}$ must attain a stationary point after a finite number of iterations. It means that the value of $C_{M C}$ cannot change after $n$th iteration, that is,

$$
\begin{align*}
& C_{M C}\left(I_{1}^{(n+1)}, \ldots, I_{j}^{(n+1)}, I_{j+1}^{(n)}, \ldots, I_{L}^{(n)} ; F^{0} ; P^{c e 0}, P^{c e 1}, N\right)  \tag{40}\\
& =C_{M C}\left(I_{1}^{(n+1)}, \ldots, I_{j-1}^{(n+1)}, I_{j}^{(n)}, \ldots, I_{L}^{(n)} ; F^{0} ; P^{c e 0}, P^{c e 1}, N\right) .
\end{align*}
$$

Using (32) and (37), we derive that $D_{j}^{(n+1)}=0$. Combining (33)-(35), we know

$$
\begin{equation*}
\left[I_{j}^{(n)}\left(Y_{i j}\right)-I_{j}^{(n+1)}\left(Y_{i j}\right)\right] G_{j}^{i}=0, \text { for } i=1, \ldots, N \tag{41}
\end{equation*}
$$

which implies either $I_{j}^{(n)}\left(Y_{i j}\right)-I_{j}^{(n+1)}\left(Y_{i j}\right)=0$, i.e., $I_{j}^{(n)}\left(Y_{i j}\right)=I_{j}^{(n+1)}\left(Y_{i j}\right)$ or $G_{j}^{i}=$ 0 , i.e., $I_{j}^{(n+1)}\left(Y_{i j}\right)=1, I_{j}^{(n)}\left(Y_{i j}\right)=0$. It follows that when $C_{M C}$ converges to a stationary
point, either $I_{j}^{(n+1)}\left(Y_{i j}\right)$ is invariant or $I_{j}^{(n+1)}\left(Y_{i j}\right)=1, I_{j}^{(n)}\left(Y_{i j}\right)=0$. Namely, $I_{j}^{(n+1)}\left(Y_{i j}\right)$ can only change from 0 to 1 at most a finite number of times. Therefore, the $I_{1}^{(n)}, I_{2}^{(n)}, \ldots, I_{L}^{(n)}$ are finitely convergent.

## 5. Extension for simultaneous search for the optimal sensor rules and fusion rule

In this section, we extend the Monte Carlo method to search for the optimal sensor rules and the optimal fusion rule simultaneously. Firstly, the necessary condition is generalized to search for the optimal sensor rules and the optimal fusion rule simultaneously. Secondly, we describe a generalized Monte Carlo Gauss-Seidel iterative algorithm. We also give the convergence of the iterative algorithm.

### 5.1 A necessary condition for the optimal sensor rules and the optimal fusion rule

Note that (15) can be rewritten as follows:

$$
\begin{align*}
& C_{M C}\left(I_{1}\left(y_{1}\right), \ldots, I_{L}\left(y_{L}\right) ; F^{0}, ; P^{c e 0},, P^{c c 1}, N\right) \\
& =c+\frac{1}{N} \sum_{i=1}^{N} \sum_{k^{\prime}=1}^{2^{L}} \sum_{k=1}^{2^{L}}\left[1-F^{0}\left(s_{k^{\prime}}\right)\right] P\left(s_{k^{\prime}} \mid s_{k}\right) \cdot P_{s_{k}}\left(\mathbf{I}\left(Y_{i}\right) \frac{\hat{L}\left(Y_{i}\right)}{g\left(Y_{i}\right)}\right.  \tag{42}\\
& =c+\frac{1}{N} \sum_{k^{\prime}=1}^{2^{L}}\left[1-F^{0}\left(s_{k^{\prime}}\right)\right] \sum_{i=1}^{N} \sum_{k=1}^{2^{L}} P\left(s_{k^{\prime}} \mid s_{k}\right) \cdot P_{s_{k}}\left(\mathbf{I}\left(Y_{i}\right) \frac{\hat{L}\left(Y_{i}\right)}{g\left(Y_{i}\right)},\right.
\end{align*}
$$

where $P_{s_{k}}\left(\mathbf{I}\left(Y_{i}\right)\right) \triangleq \prod_{j=1}^{L}\left[s_{k}(j) I_{j}\left(Y_{i j}\right)+\left(1-s_{k}(j)\right)\left(1-I_{j}\left(Y_{i j}\right)\right)\right]$ and $\mathbf{I}\left(Y_{i}\right)=$ $\left(I_{1}\left(Y_{i 1}\right), I_{2}\left(Y_{i 2}\right), \ldots, I_{L}\left(Y_{i L}\right)\right)$. Since $P_{s_{k}}\left(\mathbf{I}\left(Y_{i}\right)=1\right.$ if and only if $I_{j}=s_{k}(j)$ for all $j=$ 1, ... , $L$, (39) can be simplified as follows:

$$
\begin{align*}
& C_{M C}\left(I_{1}\left(y_{1}\right), \ldots, I_{L}\left(y_{L}\right) ; F^{0}, ; P^{c e 0},, P^{c e 1}, N\right) \\
& =c+\frac{1}{N} \sum_{k^{\prime}=1}^{2^{L}}\left[1-F^{0}\left(s_{k^{\prime}}\right)\right] \cdot \sum_{i=1}^{N} P\left(s_{k^{\prime}} \mid\left(I_{1}\left(Y_{i 1}\right), \ldots, I_{L}\left(Y_{i L}\right)\right)\right) \frac{\hat{L}\left(Y_{i}\right)}{g\left(Y_{i}\right)}, \tag{43}
\end{align*}
$$

where the terms $P_{s_{k}}\left(\mathbf{I}\left(Y_{i}\right)\right)=0$ are eliminated.
Remark 6: According to (20) and (40), the necessary condition for the optimal sensor rules is similar to Lemma 2 and the necessary condition for the optimal fusion rule is given by

$$
\begin{equation*}
F^{0}\left(s_{k^{\prime}}\right)=I\left[\sum_{i=1}^{N} P\left(s_{k^{\prime}} \mid\left(I_{1}\left(Y_{i 1}\right), \ldots, I_{L}\left(Y_{i L}\right)\right)\right) \cdot \frac{\hat{L}\left(Y_{i}\right)}{g\left(Y_{i}\right)}\right] \tag{44}
\end{equation*}
$$

for $k^{\prime}=1, \ldots, 2^{L}$. The proofs are similar to Lemma 2.

### 5.2 Generalized Monte Carlo Gauss-Seidel iterative algorithm

Based on the fixed-point type necessary condition, the generalized Monte Carlo Gauss-Seidel iterative algorithm is presented in Algorithm 2.

Remark 7: For any initial values $\left(I_{1}^{(0)}, \ldots, I_{L}^{(0)} ; F^{0(0)}\right)$, the Monte Carlo cost function $C_{M C}\left(I_{1}^{(n)}, \ldots, I_{j}^{(n)}, I_{j+1}^{(n)}, \ldots, I_{L}^{(n)} ; F^{0(n)} ; P^{c e 0}, P^{c e 1}, N\right)$ must converge to a stationary point and Algorithm 2 terminates after a finite number of iterations. The proofs are similar to those of Lemma 3 and Theorem 1.2.

Algorithm 2: Simultaneous optimization of the sensor rules and the fusion rule.

- Step 1: Generate $N$ samples: $Y_{1}, \ldots, Y_{N} \sim g(y)$, where $g(y)$ is an importance sampling density and $Y_{i}=\left[Y_{i 1}, Y_{i 2}, \ldots, Y_{i L}\right]$.
- Step 2: Initialize the $L$ sensor rules and the fusion rule, respectively, for $j=1,2, \ldots, L, i=1, \ldots, N$, and $k^{\prime}=1, \ldots, 2^{L}$,

$$
I_{j}^{(0)}\left(Y_{i j}\right)=0 / 1, F^{0(0)}\left(s_{k^{\prime}}\right)=0 / 1
$$

- Step 3: Iteratively search for the $L$ sensor rules and the fusion rule until a termination criterion in Step 4 is satisfied. The $n+1$ th iteration is given as follows: for $i=1, \ldots, N$ and $k^{\prime}=1, \ldots, 2^{L}$

$$
\begin{aligned}
& I_{1}^{(n+1)}\left(Y_{i 1}\right)=I\left[P_{11}\left(I_{2}^{(n)}\left(Y_{i 2}\right), I_{3}^{(n)}\left(Y_{i 3}\right), \ldots, I_{L}^{(n)}\left(Y_{i L}\right) ; F^{0(n)} ; P^{c e 0}, P^{c e 1}\right) \cdot \hat{L}\left(Y_{i}\right)\right], \\
& I_{2}^{(n+1)}\left(Y_{i 2}\right)=I\left[P_{21}\left(I_{1}^{(n+1)}\left(Y_{i 1}\right), I_{3}^{(n)}\left(Y_{i 3}\right), \ldots, I_{L}^{(n)}\left(Y_{i L}\right) ; F^{0(n)} ; P^{c e 0}, P^{c e 1}\right) \cdot \hat{L}\left(Y_{i}\right)\right], \\
& \ldots \ldots \\
& I_{L}^{(n+1)}\left(Y_{i L}\right)=I\left[P_{L 1}\left(I_{1}^{(n+1)}\left(Y_{i 1}\right), I_{2}^{(n+1)}\left(Y_{i 2}\right), \ldots, I_{L-1}^{(n+1)}\left(Y_{i(L-1)}\right) ; F^{0(n)} ; P^{c e 0}, P^{c e 1}\right) \cdot \hat{L}\left(Y_{i}\right)\right], \\
& \quad F^{0(n+1)}\left(s_{k^{\prime}}\right)=I\left[\sum_{i=1}^{N} P\left(s_{k^{\prime}} \mid\left(I_{1}^{(n+1)}\left(Y_{i 1}\right), I_{2}^{(n+1)}\left(Y_{i 2}\right), \ldots, I_{L}^{(n+1)}\left(Y_{i L}\right)\right)\right) \frac{\hat{L}\left(Y_{i}\right)}{g\left(Y_{i}\right)}\right] .
\end{aligned}
$$

- Step 4: For $i=1, \ldots, N$ and $k^{\prime}=1,2, \ldots, 2^{L}$, the termination criterion of the iteration process is

$$
\begin{aligned}
& I_{1}^{(n+1)}\left(Y_{i 1}\right)=I_{1}^{(n)}\left(Y_{i 1}\right), \\
& I_{2}^{(n+1)}\left(Y_{i 2}\right)=I_{2}^{(n)}\left(Y_{i 2}\right), \\
& \ldots \ldots
\end{aligned},
$$

## 6. Numerical examples

In this section, in order to evaluate the performance of Algorithms 1 and 2, we present some examples with a Gaussian signal $s$ observed in the presence of Gaussian sensor noises.

The random signal $s$ and observation noises $v_{1}, v_{2}, \ldots, v_{L}$ are as follows:

$$
\begin{equation*}
H_{0}: y_{j}=v_{j} ; \quad H_{1}: y_{j}=s+v_{j}, \quad \text { for } j=1, \ldots, L \tag{45}
\end{equation*}
$$

where $v_{1}, v_{2} \ldots, v_{L}, s$ are all mutually independent and

$$
v_{j} \sim N(0,0.6), \quad s \sim N(1,0.4), \text { for } j=1, \ldots, L
$$

Thus, given $H_{0}$ and $H_{1}$, the two conditional probability density functions are

$$
p\left(y_{1}, y_{2}, \ldots, y_{L} \mid H_{0}\right) \sim N\left(\mu_{0}, \Sigma_{0}\right), \quad p\left(y_{1}, y_{2}, \ldots, y_{L} \mid H_{1}\right) \sim N\left(\mu_{1}, \Sigma_{1}\right)
$$

where $\mu_{0}, \mu_{1}, \Sigma_{0}, \Sigma_{1}$ are easily obtained from the relationship of $s, v_{1}, v_{2}, \ldots, v_{L}$.
Assume that each sensor is required to transmit a bit through a channel with probabilities of $P_{j}^{\text {coo }}=P_{j}^{c e 1}=p$, where $p=0.05,0.15,0.3$, for $j=1,2, \ldots, L$. In the cost function (2), let the cost coefficients $C_{00}=C_{11}=0$ and $C_{10}=C_{01}=1$. The receiver operating characteristics (ROC) curves are used to evaluate the performance of the algorithms. Pf and $P d$ denote the probability of false alarm and the probability of detection, respectively.

### 6.1 Two-sensor network

We compare the Monte Carlo Gauss-Seidel iterative algorithm with the centralized algorithm and the iterative algorithm based on the Riemann sum approximation in [31] by using the receiver operating characteristics (ROC) curves.

In this case, we know $\left.\mu_{0}=[0,0]^{T}, \mu_{1}=[1,1]\right)^{T}$ and $\Sigma_{0}=[0.6,0 ; 0,0.6], \Sigma_{1}=$ $[1,0.4 ; 0.4,1]$. Some discrete values of $a$ and $b$ are used to plot ROC curves. We refer to the optimal importance sampling density $g(y) \propto\left|P_{H_{0}^{0}}(y) \hat{L}(y)\right|$ in Section 2.2 and $|\hat{L}(y)|=\left|a p\left(y \mid H_{1}\right)-b p\left(y \mid H_{0}\right)\right|$. The form is similar to the mixture-Gaussian distribution. Therefore, the importance sampling density $g(y)$ is chosen to be the mixtureGaussian distribution. The effects of choosing different $g(y)$ in terms of the performance of the Monte Carlo method were shown in [21] via numerical examples. For Algorithm 1, we take $N=200$ samples from the density $g(y)$. For the Riemann sum approximation iterative algorithm in [31], we take a discretized step-size $\Delta=0.09$, $y_{i} \in[-8,10]$, i.e., $N_{1}=N_{2}=N=200$. The ROC curves for three important fusion rules: AND, OR, and XOR rules with $p=0.05,0.15,0.3$ are plotted in Figure 2. We compare the computational time of the two algorithms with $p=0.15$ in Figure 3. Note that the analytical solution is used for the AND rule and the OR rule. Since the XOR rule is not a $K$-out-of- $L$ rule, we use Algorithm 1 to search for the sensor rules.

Some observations in Figures 2 and $\mathbf{3}$ are presented as follows:

- Given the fusion rule, the two points $(0,0)$ and $(1,1)$ may not be the beginning or ending points of the ROC curves, which is different from the case in the ideal channel cases. In addition, the larger the probability of channel errors is, the


Figure 2.
Two-sensor ROC curves with the probability of channel errors $p=0.05,0.15,0.3$.
farther away from $(0,0)$ or $(1,1)$ the ROC curves are. A possible reason is that the detection probability is not equal to 0 or 1 , even when the false alarm probability is 0 or 1 in the presence of channel errors.

- From Figure 2, when the probability of channel transmission errors increases, the decision fusion performance of different methods using the optimal sensor rules decreases.
- It can be seen in Figure 2 that the ROC curves of the new Monte Carlo approach are very close to those of the previous algorithm based on the Riemann sum approximation. However, from Figure 3, the computational time of the Monte Carlo importance sampling approximation is much less than that of the Riemann sum approximation for the three different fusion rules. It also implies that the new method can be used to deal with large-scale sensor networks.
- Note that the computational time of the AND rule and the OR rule is less than that of the XOR rule for the Monte Carlo importance sampling approximation from Figure 3. The reason is that the AND rule and the OR rule belong to the $K$-out-of- $L$ rules. The analytical form is used for the AND rule and the OR rule, therefore, the corresponding computation time is relatively lower.


Figure 3.
Two-sensor computational time as $N$ increases with the probability of channel errors $p=0.15$.

### 6.2 Ten-sensor network

We consider a larger sensor network with 10 sensors, which cannot be dealt with by the previous decision fusion algorithm based on the Riemann sum approximation due to its heavy computation requirements. For different probabilities of channel errors, the ROC curves of the AND rule, the OR rule, the 4 -out-of-10 rule, the 6 -out-of-10 rule, and Algorithm 2 are plotted in Figure 4.

Some observations in Figure 4 are presented as follows:

- The ROC curves for the ten-sensor network exhibit similar behavior as those for the two-sensor network.
- Given the fusion rule and the probability of channel errors, the decision fusion performance of the AND rule is better than the other rules for a small false alarm probability and the decision fusion performance of the OR rule is better than the other rules for a large false alarm probability. The reason may be that both of them are extreme cases of the fusion rules. For other cases, the 4-out-of-10 rule and the 6 -out-of-10 rule perform better than the AND rule and the OR rule, respectively.
- Regardless of the centralized detection algorithm, Figure 4 shows that the ROC curves generated by Algorithm 2 obtain almost the best performance for different probabilities of channel errors. It implies that a simultaneous search for the sensor rules and the fusion rule would provide better decision fusion performance.


### 6.3 One-hundred-sensor network

We consider a large-scale network with one hundred sensors. The parameter settings are similar to Section 6.2. The ROC curves of the 20 -out-of-100 rule, the 40 -out-of-100, the 50 -out-of-100, the 60 -out-of- 100 rule, and the 80 -out-of- 100 rule are plotted in Figure 5.

From Figure 5, it can be seen that the dramatically lower computational requirement of our method enables us to handle a large sensor network consisting of one hundred sensors. This is due to the fact that we have shown that there exist analytical forms of the optimal sensor rules for the $K$-out-of- $L$ rule. In addition, the decision fusion performance of different methods is improved, as the number of sensors becomes large.


Figure 4.
Ten-sensor ROC curves with the probability of channel errors $p=0.05,0.15,0.3$.


Figure 5.
One-hundred-sensor ROC curves with the probability of channel errors $p=0.05,0.15,0.3$.

## 7. Conclusion

By employing the Monte Carlo importance sampling technique, decision fusion algorithms have been provided for large-scale sensor networks with dependent observations and channel errors. The Bayesian cost function is approximated by the Monte Carlo cost function. The necessary conditions for the optimal sensor rules and the optimal fusion rule that minimize the Monte Carlo cost function have been obtained. Computationally efficient Monte Carlo Gauss-Seidel iterative algorithms have been proposed to search for the optimal sensor rules and the optimal fusion rule. These algorithms have been shown to converge after a finite number of iterations. The computational complexity of the new algorithm (i.e., $O(L N)$ ) is much less than that of the previous algorithm based on Riemann sum approximation (i.e., $O\left(L N^{L}\right)$ ). For the $K$-out-of- $L$ rule, an analytical solution has been presented for the optimal sensor rules. Simulations have demonstrated the effectiveness of Algorithms 1 and 2. Future work will include the decision fusion algorithms under the Monte Carlo framework for other networks such as tandem networks, tree networks, and sensor networks under Byzantine attack.

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[^0]:    ' In the rest of the paper, the term "sensor rules" refers to the "decision rules at the sensors."

