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# Geometric theory of non-regular separation of variables and the bi-Helmholtz equation 

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#### Abstract

The geometric theory of additive separation of variables is applied to the search for multiplicative separated solutions of the bi-Helmholtz equation. It is shown that the equation does not admit regular separation in any coordinate system in any pseudo-Riemannian space. The equation is studied in the four coordinate systems in the Euclidean plane where the Helmholtz equation and hence the bi-Helmholtz equation is separable. It is shown that the bi-Helmoltz equation admits non-trivial non-regular separation in both Cartesian and polar coordinates, while it possesses only trivial separability in parabolic and elliptic-hyperbolic coordinates. The results are applied to the study of small vibrations of a thin solid circular plate of uniform density which is governed by the bi-Helmholtz equation.


## 1 Introduction

The solution of boundary value problems for the partial differential equations of mathematical physics by the method of separation of variables is an effective method that has been employed for almost two centuries [10, 9]. The method, which assumes that the solution is a product of functions each one of a single independent variable (product ansatz), reduces the partial differential equation to a set of ordinary differential equations with a corresponding set of separation constants. The application of the method to the various equations of mathematical physics leads to certain ordinary differential equations the solutions of which are studied in their own right in special function theory. The special functions defined by an equation arising from a given separable coordinate system define function spaces that yield the solution of the given boundary problem under general boundary values. The theory of this method for a general second order linear PDE with variable coefficients under the product ansatz is well developed [1, 7, 8] However, the theory for equations of higher order has only more recently been studied [6] where general separation is considered and where the concept of non-regular separation may be utilized. In this article a geometric formulation of this theory [4] is applied to the bi-Helmholtz equation which arises in the study of small vibrations of a thin, solid plate [11, 5].

The plan of the paper is as follows. In Section 2 the geometric theory of separation of variables is reviewed and the concepts of regular and non-regular separation defined. In Section 3 the theory of the previous section is applied to the bi-Helmholtz equation. The main result of this section is the proof of Theorem 1 which states that regular multiplicative separation for the bi-Helmholtz equation (7) on any Riemannian or pseudo-Riemannian n-dimensional manifold does not occur in any system of coordinates. Conditions for non-regular separation to occur are given in Proposition 2 and the concept of non-trivial non-regular separation is introduced. Section 4 contains examples in $\mathbb{E}^{2}$ of coordinate systems where the bi-Helmholtz equation admits trivial and non-trivial non-regular separation. In Section 5 the results of the preceding section are applied to the analysis of the small vibrations of a thin solid circular plate which reproduce the classical results obtained by Rayleigh [11]. The Conclusion is given in Section 6.

## 2 Geometric theory of separation of variables

One of the most known ansatz for solving a PDE depending on $n$ independent variables $\left(q^{i}\right)$ is the additive separation, that is the search of those solutions $u\left(q^{i}\right)$ which are written as sum of functions depending on a single variable (in the given coordinate system $\left.\left(q^{i}\right)\right)$ :

$$
u=\sum_{i=1}^{n} S_{(i)}\left(q^{i}\right)
$$

Moreover, we do not look for a single separated solution but we want to determine a family (as big as possible) of separated solutions

$$
u=\sum_{i=1}^{n} S_{(i)}\left(q^{i}, c_{\alpha}\right)
$$

where the real constant parameters $c_{\alpha}$ satisfy a suitable completeness condition (see below). When the method of SoV works, it is possible to split the PDE into $n$ ODEs of order $l$ which involve only one of the functions $S_{(i)}$; these ODEs are known as separated equations. Separability of a PDE strongly depends on the choice of the independent variables and is destroyed by a general transformation of the $\left(q^{i}\right)$. However, we want to provide a geometrical interpretation of what we mean by solving a PDE through the ansatz of separation of variables. We will distinguish between two possible types of families of separated solutions: the first type depends on the maximal number of parameters (we refer to this case as "free" [1] or "regular" [6] separation of variables) and the second type depending on less parameters. We refer to this case as "non regular separation" [6, 3].

We consider the separability of the $l$-th order PDE

$$
\begin{equation*}
\mathcal{H}\left(q^{i}, u, u_{i}, \ldots, u_{i}^{(l)}\right)=h \quad(h \in \mathbb{R}) \tag{1}
\end{equation*}
$$

in the coordinates $\left(q^{i}\right)$ on the $n$-dimensional manifold $Q$. Here and in the following we denote by $u=u\left(q^{i}\right)$ the unknown function and set

$$
u_{i}=\frac{\partial u}{\partial q^{i}}, \quad u_{i}^{(2)}=u_{i i}=\frac{\partial^{2} u}{\left(\partial q^{i}\right)^{2}}, \quad \ldots \quad u_{i}^{(l)}=\frac{\partial^{l} u}{\left(\partial q^{i}\right)^{l}} .
$$

For the sake of simplicity, we assume that the maximal order of derivatives involved in $H$ is the same for each index $i=1, \ldots, n$. Let $Z$ be the $(n l+1)$ dimensional space of the dependent variable and its separated derivatives: coordinates on $Z$ are given by $\left(u, u_{i}, u_{i}^{(2)}, \ldots, u_{i}^{(l)}\right)$. We consider the trivial bundle over $Q, M=Q \times Z$. In [1] free separation is defined as the existence of an additively separated solution $u$ of (1), depending on $n l+1$ parameters $\left(c_{A}\right)$, satisfying the completeness condition

$$
\operatorname{rank}\left[\left.\frac{\partial u}{\partial c_{A}}\left|\frac{\partial u_{i}}{\partial c_{A}}\right| \ldots \right\rvert\, \frac{\partial u_{i}^{(l)}}{\partial c_{A}}\right]=n l+1
$$

Free separation occurs if and only if the $n$ vector fields of the form

$$
\begin{equation*}
D_{i}=\frac{\partial}{\partial q^{i}}+u_{i} \frac{\partial}{\partial u}+u_{i}^{(2)} \frac{\partial}{\partial u_{i}}+\ldots+u_{i}^{(l)} \frac{\partial}{\partial u_{i}^{(l-1)}}+R_{i} \frac{\partial}{\partial u_{i}^{(l)}}, \tag{2}
\end{equation*}
$$

where $R_{i}\left(q^{j}, u, u_{j}, \ldots, u_{j}^{(l)}\right)$ are functions on $M$, are commuting symmetries of (1) namely, the conditions

$$
\begin{gather*}
D_{i} \mathcal{H}=0  \tag{3}\\
{\left[D_{i}, D_{j}\right]=0} \tag{4}
\end{gather*}
$$

are satisfied (see [1]). Eqs (3) determine the functions $R_{i}$, while (4) are equivalent to

$$
\begin{equation*}
D_{i} R_{j}=0 \quad(i \neq j) \tag{5}
\end{equation*}
$$

By expanding conditions (5), we get exactly the conditions given in [6] for regular separation. Thus, free separation corresponds to regular separation.

Remark 1. The geometrical interpretation of the free separation is summarized in the following items (see [1]):

- $\Delta=\operatorname{span}\left(D_{i}\right)$ is an integrable distribution of rank $n$ on $M$.
- The foliation of $n$-dimensional integral manifolds of $\Delta$ is described by a complete separated solution of $\mathcal{H}=h$ :

$$
\left\{\begin{array}{l}
u=S=\sum_{i=1}^{n} S_{(i)}\left(q^{i}, c_{\alpha}\right) \\
u_{i}=S_{(i)}^{\prime} \\
u_{i i}=S_{(i)}^{\prime \prime} \\
\cdots
\end{array} \quad \text { with } \partial_{i} S=S_{(i)}^{\prime}, \partial_{i j}^{2} S=\delta_{i j} S_{(i)}^{\prime \prime}, \ldots\right.
$$

- Free (regular) separated solutions depend on $n l+1$ parameters.
- Completeness means that for any point $P_{0}$ with coordinates $\left(q_{0}^{i}\right)$ there is a separated solution of $\mathcal{H}=h$ for each choice of the value of $u$ and its $n l$ derivatives $u_{i}, u_{i}^{2}, \ldots, u_{i}^{l}$ at $P_{0}$.

Remark 2. If $\mathcal{H}$ does not depend on $u$, then $u$ is defined up to an additive constant and the relevant constants appearing in the separated solutions are in fact $n l$ (the remaining one being the finial additive constant). In this case we can eliminate the variable $u$ from the space $Z$ and (2) reduce to

$$
\begin{equation*}
D_{i}=\frac{\partial}{\partial q^{i}}+u_{i}^{(2)} \frac{\partial}{\partial u_{i}}+\ldots+u_{i}^{(l)} \frac{\partial}{\partial u_{i}^{(l-1)}}+R_{i} \frac{\partial}{\partial u_{i}^{(l)}} . \tag{6}
\end{equation*}
$$

Moreover, even the separated equations do not depend on $S_{(i)}$. Thus, each of $S_{(i)}$ is defined by its ODE up to an additive constant which can be disregarded, since it only affect the inessential global additive constant of $u$. Hence, each $l$-th order ODE actually contribute with $l-1$ essential integration constants. This means that the $n l$ constants which are involved in the regular separated solutions splits into $n$ separation constants plus $n l-n=n(l-1)$ integrating constants (see Remarks 4.8 and 4.9 in [1] for the case of the Schrödinger equation).

We consider now the second type of separation, occurring when condition (5) is not identically satisfied. This case was firstly considered by Kalnins and Miller in [6] and called non-regular separation. The authors state that also in this
case separable solutions still may exist, but they will depend on less than $n l+1$ parameters.

However, in [6] it is not specified under what conditions separable solutions will exist and how to determine the number of the parameters involved.

In [3] one of us gives a geometric interpretation of the situation that naturally leads to an effective definition of non-regular separation.

Let $N$ be a submanifold of $M=Q \times Z$ locally described by the $r$ equations

$$
f_{a}=0 \quad(a=1, \ldots, r)
$$

If the vectors $D_{i}$ commute on $N$, that is

$$
\left.D_{i} R_{j}\right|_{N}=0 \quad(i \neq j, i, j=1, \ldots, n)
$$

and the vector fields $D_{i}$ are tangent to $N$, that is

$$
\left.D_{i} f_{a}\right|_{N}=0 \quad(i=1, \ldots, n, a=1, \ldots r)
$$

then we can restrict the distribution $\Delta$ generated by the $D_{i}$ to an involutive distribution on $N$. Thus, on $N$ we have a complete separated solution of $\mathcal{H}=h$ depending on $n l+1-r$ parameters $\left(c_{\alpha}\right)$, which can be considered as a constrained separated solution on $M$. Hence, according to [3] we have the following

Definition 1. The PDE $\mathcal{H}=h$ admits a non-regular or constrained additive separation on a submanifold $N$ of $M$ defined by the $r$ equations $f_{a}=0$, if

1. $u=\sum_{i} S^{(i)}\left(q^{i}, c_{\alpha}\right)$ is a solution of $\mathcal{H}=h$;
2. $u$ depends on $n l+1-r$ parameters $\left(c_{\alpha}\right)$ satisfying the completeness conditions

$$
\operatorname{rank}\left[\left.\frac{\partial u}{\partial c_{\alpha}}\left|\frac{\partial u_{i}}{\partial c_{\alpha}}\right| \ldots \right\rvert\, \frac{\partial u_{i}^{(l)}}{\partial c_{A}}\right]=n l+1-r ;
$$

3. $u$ and its derivatives satisfy $f_{\alpha}\left(q^{j}, u, u_{i} \ldots\right)=0$ for all (admissible) values of the parameters $\left(c_{\alpha}\right)$.

From the above discussion we get the following criterion for the constrained separation

Proposition 1. In a given coordinate system $\left(q^{i}\right)$ equation (1) admits a constrained separation on the submanifold $N$ defined by equations $f_{a}=0$ if and only if the vector fields $D_{i}$ (2) are symmetries of (11), tangent to $N$ and commute on $N$, that is

$$
D_{i} \mathcal{H}=0,\left.\quad D_{i} R_{j}\right|_{N}=0,\left.\quad D_{i} f_{a}\right|_{N}=0 \quad(i \neq j=1, \ldots, n, a=1, \ldots, r)
$$

Remark 3. If $D_{i} R_{j}$ are everywhere different from zero, then equation (1) has no additive separable solutions. A possible choice for $N$ is the set of points satisfying equations $D_{i} R_{j}=0$, but in many cases this set is not a well-defined manifold or the vectors $D_{i}$ are not tangent to it.

Remark 4. In analogy with Remark [1, the geometric interpretation of the constrained separation on a submanifold $N$ is sketched in the following items:

- $\Delta_{S}=\operatorname{span}\left(D_{i}\right)$ is an integrable $n$-dimensional distribution on $N$ whose integral manifolds are described by a constrained complete separated solution of $\mathcal{H}=h$.
- constrained separated solutions depend on $n l+1-r$ parameters.
- completeness means that for any $P_{0}=\left(q_{0}^{i}\right)$ there is a separated solution of $\mathcal{H}=h$ for each choice of the value of $u$ and its $n l$ separated derivatives at $P_{0}$ satisfying $f_{\alpha}=0$ (i.e., the initial condition belongs to $N$ ).


## 3 Application to the bi-Helmholtz equation

In order to apply the theory recalled in the previous section to the multiplicative separation of the bi-Helmholtz equation, we should transform the unknown function $\psi$ to its logarithm $u=\log \psi$ and then we should write down the simplified equation for separated solutions. First of of all, we need to express the bi-Helmholtz equation

$$
\begin{equation*}
\Delta^{2} \psi=\lambda \psi, \quad \lambda \in \mathbb{R} \tag{7}
\end{equation*}
$$

where $\Delta$ is the Laplace-Beltrami operator on a $n$-dimensional Riemannian or pseudo-Riemannian manifold $(Q, \mathbf{g})$ with respect to a coordinate system $\left(q^{i}\right)$ in which the metric has contravariant components $\left(g^{i j}\right)$. For the moment we do not require any special property to the coordinates (such as orthogonality, or specifying of the dimension $n$ ). It is well known that

$$
\begin{equation*}
\Delta \psi=g^{i j} \partial_{i} \partial_{j} \psi-\Gamma^{i} \partial_{i} \psi \tag{8}
\end{equation*}
$$

where $\partial_{i}$ is the partial derivative with respect to $q^{i}$ and

$$
\Gamma^{h}=g^{i j} \Gamma_{i j}^{h}=\frac{1}{2} g^{i j} g^{h k}\left(\partial_{i} g_{j k}+\partial_{j} g_{i k}-\partial_{k} g_{i j}\right)
$$

are the contracted Christoffel symbols (see [1]).
If not explicitly written, Einstein summation notation is sytematically employed: repeated upper and lower indices means summation.

By applying twice the Laplace operator we get

$$
\begin{align*}
\Delta^{2} \psi= & g^{i j} \partial_{i j}\left(g^{h k} \partial_{h k} \psi-\Gamma^{h} \partial_{h} \psi\right)-\Gamma^{i} \partial_{i}\left(g^{h k} \partial_{h k} \psi-\Gamma^{h} \partial_{h} \psi\right)= \\
= & g^{i j} g^{h k} \partial_{i j h k} \psi+2\left(g^{i j} \partial_{j} g^{h k}-g^{h k} \Gamma^{i}\right) \partial_{i h k} \psi+ \\
& \left(g^{i j} \partial_{i j} g^{h k}-2 g^{j k} \partial_{j} \Gamma^{h}-\Gamma^{i} \partial_{i} g^{h k}+\Gamma^{h} \Gamma^{k}\right) \partial_{h k} \psi+  \tag{9}\\
& \left(-g^{i j} \partial_{i j} \Gamma^{h}+\Gamma^{i} \partial_{i} \Gamma^{h}\right) \partial_{h} \psi,
\end{align*}
$$

where

$$
\begin{equation*}
\partial_{i_{1} \ldots i_{p}}:=\partial_{i_{1}} \ldots \partial_{i_{p}} . \tag{10}
\end{equation*}
$$

In order to simplify notation we can write the bi-Laplacian as

$$
\begin{equation*}
\Delta^{2} \psi=A^{i j k l} \partial_{i j k l} \psi+B^{i j k} \partial_{i j k} \psi+C^{i j} \partial_{i j} \psi+D^{i} \partial_{i} \psi \tag{11}
\end{equation*}
$$

where

$$
\begin{align*}
& A^{i j k l}=g^{(i j} g^{k l)}  \tag{12}\\
& B^{i j k}=2\left(g^{h(i} \partial_{h} g^{j k)}-g^{(i j} \Gamma^{k)}\right)  \tag{13}\\
& C^{i j}=g^{k l} \partial_{k l} g^{i j}-2 g^{k(i} \partial_{k} \Gamma^{j)}-\Gamma^{k} \partial_{k} g^{i j}+\Gamma^{i} \Gamma^{j}  \tag{14}\\
& D^{i}=-g^{j k} \partial_{j k} \Gamma^{i}+\Gamma^{j} \partial_{j} \Gamma^{i}=-\Delta \Gamma^{i} \tag{15}
\end{align*}
$$

where $(\cdots)$ indicates symmetrization of the indices.
In order to pass from a multiplicative separated solution to an additive separated solution we perform the change of unknown

$$
\psi=e^{u}, \quad u=\log \psi
$$

We denote the partial derivatives of $u$ by $u_{i}, u_{i j}$ etc. The link between partial derivatives of $\psi$ and $u$ is given by

$$
\begin{aligned}
& \partial_{i} \psi=e^{u} u_{i} \\
& \partial_{i j} \psi=e^{u}\left(u_{i j}+u_{i} u_{j}\right) \\
& \partial_{i j k} \psi=e^{u}\left(u_{i j k}+3 u_{(i} u_{j k)}+u_{i} u_{j} u_{k}\right) \\
& \partial_{i j k l} \psi=e^{u}\left(u_{i j k l}+4 u_{(i} u_{j k l)}+3 u_{(i j} u_{k l)}+6 u_{(i} u_{j} u_{k l)}+u_{i} u_{j} u_{k} u_{l}\right)
\end{aligned}
$$

Hence, the equation to which we want to apply the geometric theory of SoV is

$$
H\left(q^{i}, u_{i}, u_{i j}, \ldots, u_{i j k l}\right)=\lambda, \quad \lambda \in \mathbb{R}
$$

where

$$
\begin{aligned}
H= & A^{i j k l} u_{i j k l}+\left(B^{i j k}+4 A^{i j k l} u_{l}\right) u_{i j k}+\left(C^{i j}+3 B^{i j k} u_{k}\right) u_{i j}+ \\
& 3 A^{i j k l} u_{i j} u_{k l}+6 A^{i j k l} u_{i} u_{j} u_{k l}+ \\
& A^{i j k l} u_{i} u_{j} u_{k} u_{l}+B^{i j k} u_{i} u_{j} u_{k}+C^{i j} u_{i} u_{j}+D^{i} u_{i},
\end{aligned}
$$

If we are interested only in separated solutions

$$
u=\sum_{i} S_{(i)}\left(q^{i}, c_{A}\right), \quad A=1, \ldots, 4 n
$$

where, since the PDE does not depend explicitly on $u$, the constants $\left(c_{A}\right)$ are at most $4 n$, we replace $H$ by

$$
\begin{align*}
H_{s}= & \left(g^{i i}\right)^{2} u_{i}^{(4)}+\left(4 g^{i i} g^{i j} u_{j}+B^{i i i}\right) u_{i}^{(3)}+\left(g^{i i} g^{j j}+2\left(g^{i j}\right)^{2}\right) u_{i}^{(2)} u_{j}^{(2)}+ \\
& \left(2\left(g^{i i} g^{h j}+2 g^{i j} g^{i h}\right) u_{j} u_{h}+\left(B^{i i j}+B^{i j i}+B^{j i i}\right) u_{j}+C^{i i}\right) u_{i}^{(2)}+  \tag{16}\\
& g^{i j} g^{h k} u_{i} u_{j} u_{h} u_{k}+B^{i h k} u_{i} u_{h} u_{k}+C^{i j} u_{i} u_{j}+D^{i} u_{i},
\end{align*}
$$

where $u_{i}^{(s)}$ is the $s$-th partial derivative of $u$ w.r.t. $q^{i}$.
The function (16) is a 4 -th degree polynomial in the derivatives of $u$ to the 4 -th order; by observing that the terms $B^{i j h}, C^{i j}, D^{i}$ contains first, second, third order derivatives of the metric tensor respectively, we can say that $H_{s}$ is homogeneous: the sum of the degree and the order of the derivatives it is the same (four) in all terms.

For the further computation it is useful to write down the first partial derivatives of $H_{s}$

$$
\begin{align*}
\frac{\partial H_{s}}{\partial u_{i}^{(4)}}= & \left(g^{i i}\right)^{2}  \tag{17}\\
\frac{\partial H_{s}}{\partial u_{i}^{(3)}}= & B^{i i i}+4 g^{i i} g^{i j} u_{j}  \tag{18}\\
\frac{\partial H_{s}}{\partial u_{i}^{(2)}}= & 2\left(g^{i i} g^{j j}+2\left(g^{i j}\right)^{2}\right) u_{j}^{(2)}+2\left(g^{i i} g^{h j}+2 g^{i j} g^{i h}\right) u_{j} u_{h}+  \tag{19}\\
& \left(B^{i i j}+B^{i j i}+B^{j i i}\right) u_{j}+C^{i i} \\
\frac{\partial H_{s}}{\partial u_{i}}= & 4 g^{j j} g^{i j} u_{j}^{(3)}+u_{j}^{(2)}\left(B^{j j i}+B^{j i j}+B^{i j j}+4\left(g^{j j} g^{h i}+2 g^{j i} g^{j h}\right) u_{h}\right)+(20)  \tag{20}\\
& D^{i}+\left(C^{i j}+C^{j i}\right) u_{j}+\left(B^{i j h}+B^{j h i}+B^{h i j}\right) u_{j} u_{h}+4 g^{i j} g^{h k} u_{j} u_{h} u_{k}
\end{align*}
$$

We compute the vector fields associated with the separation, which are of the form (6):

$$
D_{i}=\partial_{i}+u_{i}^{(2)} \frac{\partial}{\partial u_{i}}+u_{i}^{(3)} \frac{\partial}{\partial u_{i}^{(2)}}+u_{i}^{(4)} \frac{\partial}{\partial u_{i}^{(3)}}+R_{i} \frac{\partial}{\partial u_{i}^{(4)}},
$$

where the index $i$ is not summed and the functions $R_{i}\left(q^{h}, u_{h}, \ldots u_{h}^{(4)}\right)$ are determined by the condition

$$
D_{i}\left(H_{s}\right)=0 .
$$

Hence, under the technical assumption that $g^{i i} \neq 0$, we have

$$
R_{i}=-\left(\frac{\partial H_{s}}{\partial u_{i}^{(4)}}\right)^{-1}\left(\partial_{i} H_{s}+u_{i}^{(2)} \frac{\partial H_{s}}{\partial u_{i}}+u_{i}^{(3)} \frac{\partial H_{s}}{\partial u_{i}^{(2)}}+u_{i}^{(4)} \frac{\partial H_{s}}{\partial u_{i}^{(3)}}\right)=
$$

$$
=-\frac{1}{\left(g^{i i}\right)^{2}}\left(\partial_{i} H_{s}+u_{i}^{(2)} \frac{\partial H_{s}}{\partial u_{i}}+u_{i}^{(3)} \frac{\partial H_{s}}{\partial u_{i}^{(2)}}+u_{i}^{(4)} \frac{\partial H_{s}}{\partial u_{i}^{(3)}}\right)
$$

The condition for free (or regular) separation is that

$$
D_{i} R_{j}=0 \quad \forall i \neq j
$$

that is

$$
\partial_{i} R_{j}+u_{i}^{(2)} \frac{\partial R_{j}}{\partial u_{i}}+u_{i}^{(3)} \frac{\partial R_{j}}{\partial u_{i}^{(2)}}+u_{i}^{(4)} \frac{\partial R_{i}}{\partial u_{i}^{(3)}}+R_{i} \frac{\partial R_{j}}{\partial u_{i}^{(4)}}=0 .
$$

We look for terms containing $u_{h}^{(3)}$ or $u_{h}^{(4)}$ and such that the sum of the orders of the derivatives of $u$ is six i.e., the terms containing $u_{h}^{(4)} u_{k}^{(2)}$ or $u_{h}^{(3)} u_{k}^{(3)}$ (for $k \neq h$ ) and with the components of the metric tensor not derived. Since we have for $i \neq j$

$$
\begin{gathered}
\frac{\partial R_{j}}{\partial u_{i}^{(4)}}=-\frac{1}{\left(g^{j j}\right)^{2}}\left(\frac{\partial^{2} H_{s}}{\partial q^{j} \partial u_{i}^{(4)}}+u_{j}^{(2)} \frac{\partial^{2} H_{s}}{\partial u_{i}^{(4)} \partial u_{j}}+u_{j}^{(3)} \frac{\partial^{2} H_{s}}{\partial u_{i}^{(4)} \partial u_{j}^{(2)}}+u_{j}^{(4)} \frac{\partial^{2} H_{s}}{\partial u_{j}^{(3)} \partial u_{i}^{(4)}}\right) \\
=-\frac{\partial_{j}\left(g^{i i}\right)^{2}}{\left(g^{j j}\right)^{2}},
\end{gathered}
$$

the addendum $R_{i} \frac{\partial R_{j}}{\partial u_{i}^{(4)}}$ does not contain such a term, as well as $\partial_{i} R_{j}$. Moreover, we have

$$
\begin{gathered}
\frac{\partial R_{j}}{\partial u_{i}^{(3)}}=-\frac{1}{\left(g^{j j}\right)^{2}} \\
\left(\frac{\partial^{2} H_{s}}{\partial q^{j} \partial u_{i}^{(3)}}+u_{j}^{(2)} \frac{\partial^{2} H_{s}}{\partial u_{i}^{(3)} \partial u_{j}}+u_{j}^{(3)} \frac{\partial^{2} H_{s}}{\partial u_{i}^{(3)} \partial u_{j}^{(2)}}+u_{j}^{(4)} \frac{\partial^{2} H_{s}}{\partial u_{j}^{(3)} \partial u_{i}^{(3)}}\right) \\
=-\frac{1}{\left(g^{j j}\right)^{2}}\left(\partial_{j} B^{i i i}+4 \partial_{j}\left(g^{i i} g^{i h}\right) u_{h}+4 u_{j}^{(2)} g^{i i} g^{i j}\right)
\end{gathered}
$$

Thus, in the addendum $u_{i}^{(4)} \frac{\partial R_{j}}{\partial u_{i}^{(3)}}$, the function $D_{i} R_{j}$ contains the term

$$
\begin{equation*}
-4 \frac{g^{i i} g^{i j}}{\left(g^{j j}\right)^{2}} u_{j}^{(2)} u_{i}^{(4)} \tag{21}
\end{equation*}
$$

For the derivative of $R_{j}$ w.r.t. $u_{i}^{(2)}$ we have

$$
\begin{aligned}
\frac{\partial R_{j}}{\partial u_{i}^{(2)}}= & -\frac{1}{\left(g^{j j}\right)^{2}}\left(\frac{\partial^{2} H_{s}}{\partial q^{j} \partial u_{i}^{(2)}}+u_{j}^{(2)} \frac{\partial^{2} H_{s}}{\partial u_{i}^{(2)} \partial u_{j}}+u_{j}^{(3)} \frac{\partial^{2} H_{s}}{\partial u_{i}^{(2)} \partial u_{j}^{(2)}}+u_{j}^{(4)} \frac{\partial^{2} H_{s}}{\partial u_{j}^{(3)} \partial u_{i}^{(2)}}\right)= \\
& =-\frac{1}{\left(g^{j j}\right)^{2}}\left(\frac{\partial^{2} H_{s}}{\partial q^{j} \partial u_{i}^{(2)}}+u_{j}^{(2)} \frac{\partial^{2} H_{s}}{\partial u_{i}^{(2)} \partial u_{j}}+2\left(g^{i i} g^{j j}+2\left(g^{i j}\right)^{2}\right) u_{j}^{(3)}\right)
\end{aligned}
$$

Therefore, in $u_{i}^{(3)} \frac{\partial R_{j}}{\partial u_{i}^{(2)}}$, the function $D_{i} R_{j}$ contains the term

$$
\begin{equation*}
\frac{-2}{\left(g^{j j}\right)^{2}}\left(g^{i i} g^{j j}+2\left(g^{i j}\right)^{2}\right) u_{j}^{(3)} u_{i}^{(3)} \tag{22}
\end{equation*}
$$

and, since $\frac{\partial H_{s}}{\partial u_{i}^{(2)}}$ does not depend on $u_{h}^{(3)}$ and $u_{h}^{(4)}$, we disregard other addenda of $\frac{\partial R_{j}}{\partial u_{i}^{(2)}}$. The last term to consider is

$$
\begin{aligned}
& \frac{\partial R_{j}}{\partial u_{i}}=-\frac{1}{\left(g^{j j}\right)^{2}}\left(\frac{\partial^{2} H_{s}}{\partial q^{j} \partial u_{i}}+u_{j}^{(2)} \frac{\partial^{2} H_{s}}{\partial u_{i} \partial u_{j}}+u_{j}^{(3)} \frac{\partial^{2} H_{s}}{\partial u_{i} \partial u_{j}^{(2)}}+u_{j}^{(4)} \frac{\partial^{2} H_{s}}{\partial u_{j}^{(3)} \partial u_{i}}\right) \\
= & -\frac{1}{\left(g^{j j}\right)^{2}}\left(4 g^{j j} g^{i j} u_{j}^{(4)}+\text { terms with derivatives of } u \text { of order less than } 4^{t h}\right)
\end{aligned}
$$

Thus $D_{i} R_{j}$ contains the term

$$
\begin{equation*}
-4 \frac{g^{i j}}{g^{j j}} u_{j}^{(4)} u_{i}^{(2)} \tag{23}
\end{equation*}
$$

and no other terms in $u_{h}^{(4)} u_{k}^{(2)}$ or $u_{h}^{(3)} u_{k}^{(3)}$ are present in $D_{i} R_{j}$ except for (21,|22],23). It is easy to see that the three terms cannot be all zero since $g^{i i} \neq 0$, for all $i$. Hence, we have proved that

Theorem 1. Regular multiplicative separation for the bi-Helmholtz equation (7) on any Riemannian or pseudo-Riemannian n-dimensional manifold does not occur in any system of coordinates.

Remark 5. This means that we cannot find a family of separated solutions that depends on $4 n$ parameters and such that the values of

$$
\begin{equation*}
u_{i}, \quad u_{i}^{(2)}, \quad u_{i}^{(3)}, \quad u_{i}^{(4)} \tag{24}
\end{equation*}
$$

can be fixed in an arbitrary way. Going back to the multiplicatively separated solutions $\prod_{i} \psi_{i}$ of (17), the impossibility that regular separation occurs can be interpreted as follows: the one-to-one relationship existing between (24) and

$$
\frac{\psi_{i}^{\prime}}{\psi_{i}} \quad \frac{\psi_{i}^{\prime \prime}}{\psi_{i}} \quad \frac{\psi_{i}^{\prime \prime \prime}}{\psi_{i}} \quad \frac{\psi_{i}^{\prime \prime \prime \prime}}{\psi_{i}}
$$

(see [1] for the explicit link till order 2) implies that it is not possible to assign these values in an arbitrary way at a point $q_{0}$ and finding a separated solution satisfying these initial conditions.

However, smaller families of separated solutions may exist. In this case the theory of non-regular separation (see [3]) may be usefully applied.

Proposition 2. In any coordinate system ( $q^{i}$ ) allowing regular separation for the Helmholtz equation, there exists a submanifold $N$ of $M$ with dimension at least $2 n$ where bi-Helmholtz admits non-regular separation.

The statement follows form the fact that every solution of the Helmholtz equation is also a solution of the bi-Helmholtz equation. Hence, the $2 n$-parameter families of the multiplicatively separated solutions of the Helmholtz equation are multiplicatively separated solution of (7). Non-regular separation separation is said to be non-trivial if there exists a bigger family of separable solutions for the bi-Helmholtz equation than that for the Helmholtz equation. Otherwise, the separation is said to be trivial. A detailed analysis of some examples shows that non-trivial non-regular separation is possible.

## 4 Examples

In this section we provide some explicit computations of the manifolds $N$ allowing non-regular separation of variables for the bi-Helmholtz equation. For sake of simplicity, we restrict ourselves to the Euclidean plane. Moreover, in order to be sure that there exists a $N$ where non-regular separation occurs, we shall consider only coordinates that allow separation for the Helmholtz equation. Several situations will be described. In two of the four possible separable coordinate systems (Cartesian and polar) we show that non-trivial non-regular separation of the bi-Helmholtz equation is possible (sections 4.1 and 4.2). One the other hand, if we use parabolic or elliptic-hyperbolic coordinates, the only possibility is trivial separation.

### 4.1 Cartesian coordinates on the plane

Let us examine the case of bi-Helmholtz equation on the plane in Cartesian coordinates $\left(q^{1}, q^{2}\right)=(x, y)$. The only separability condition $D_{1} R_{2}=0$ becomes

$$
\begin{equation*}
-\left(2 u_{2}^{(2)} u_{2}+u_{2}^{(3)}\right)\left(2 u_{1}^{(2)} u_{1}+u_{1}^{(3)}\right)=0 \tag{25}
\end{equation*}
$$

which is satisfied if one of the factors vanishes. However, this condition alone does not define a submanifold such that $D_{1}$ and $D_{2}$ are tangent to it: for example

$$
\begin{equation*}
D_{1}\left(2 u_{1}^{(2)} u_{1}+u_{1}^{(3)}\right)=2\left(u_{1}^{(2)}\right)^{2}+2 u_{1}^{(3)} u_{1}+u_{1}^{(4)} \tag{26}
\end{equation*}
$$

Thus this condition must also be added in order to define a manifold where the non-regular separation could occur. Let us call $N$ the submanifold defined by the equations

$$
\begin{equation*}
f_{1}=u_{1}^{(3)}+2 u_{1}^{(2)} u_{1}=0, \quad f_{2}=u_{1}^{(4)}+2\left(u_{1}^{(2)}\right)^{2}+2 u_{1}^{(3)} u_{1}=0 \tag{27}
\end{equation*}
$$

It is easy to check that $\left.D_{i} f_{a}\right|_{N}=0(i=1,2 a=1,2)$ and that $\left.D_{i} R_{j}\right|_{N}=$ 0 . Hence, on the 6 -dimensional manifold $N$ non-regular separation occurs in Cartesian coordinates. The four derivatives of $S_{2}(y)\left(u_{2}, u_{2}^{(2)}, u_{2}^{(3)}, u_{2}^{(4)}\right)$ can be arbitrarily assigned at an initial point, while only two of the derivatives of $S_{1}(x)$ $\left(u_{1}, u_{1}^{(2)}\right)$ are free and the remaining ones are determined by the equations $f_{a}$ defining $N$.

Since the equations for $N$ only involve the dependence of $u$ on $q^{1}=x$, they can be seen as separated equations for the function $S_{1}\left(q^{1}\right)$ : in particular $f_{1}$ and $f_{2}$ (its differential consequence) imply that $S_{1}$ has to satisfy

$$
\begin{equation*}
\frac{d^{2}}{d x^{2}} S_{1}+\left(\frac{d}{d x} S_{1}\right)^{2}=c_{1}, \quad c_{1} \in \mathbb{R} \tag{28}
\end{equation*}
$$

It is interesting to remark that (28) means precisely that $S_{1}$ is a separated solution of the Laplace equation. Indeed, by inserting the function $\psi_{1}(x)=e^{S_{1}}$ in (28), we get

$$
\psi_{1}^{\prime \prime}=c_{1} \psi_{1} \quad\left(c_{1} \in \mathbb{R}\right)
$$

However, the family of the separated solutions is bigger than the separated solutions of the Helmholtz equation, since for $S_{2}$ we can get functions which are not solutions of the Helmholtz equation but only of the bi-Helmholtz one. Indeed, by (28) the separated equation for $H_{s}$ written in Cartesian coordinates takes the form

$$
\begin{equation*}
u_{2}^{(4)}+3\left(u_{2}^{(2)}\right)^{2}+4 u_{2} u_{2}^{(3)}+6\left(u_{2}\right)^{2} u_{2}^{(2)}+\left(u_{2}\right)^{4}+2 c_{1}\left(u_{2}^{(2)}+\left(u_{2}\right)^{2}\right)+c_{1}^{2}-\lambda=0 \tag{29}
\end{equation*}
$$

where $c_{1}=u_{1}^{(2)}+\left(u_{1}\right)^{2}$ which is constant on the constraint surface $N$. A simpler form for (29) is given in terms of $\psi_{2}=\ln \left(S_{2}\right)$ : with respect to this new independent variable the ODE becomes the linear 4 -th order equation

$$
\left(\psi_{2}\right)^{(4)}+2 c_{1} \psi_{2}^{\prime \prime}+\left(c_{1}^{2}-\lambda\right) \psi_{2}=0
$$

### 4.2 Polar coordinates

Let us examine the case of bi-Helmholtz equation on the plane in polar coordinates $\left(q^{1}, q^{2}\right)=(r, \theta)$. The situation is similar, but is no longer symmetric in the two variables because we pass from two ignorable coordinates to one ignorable coordinate.

$$
\begin{equation*}
\left(u_{1}^{(3)} r^{2}+2 r\left(u_{1}\right)^{2}+u_{1}^{(2)} r-u_{1}+2 u_{1}^{(2)} u_{1} r^{2}\right)\left(2 u_{2}^{(2)} u_{2}+u_{2}^{(3)}\right)=0 \tag{30}
\end{equation*}
$$

As in the Cartesian case, if we consider the condition on $\theta$-depending function, on the six-dimensional submanifold $N_{\text {pol }}$ defined by

$$
\begin{equation*}
f_{1}=2 u_{2}^{(2)} u_{2}+u_{2}^{(3)}=0, \quad f_{2}=u_{2}^{(4)}+2\left(u_{2}^{(2)}\right)^{2}+2 u_{2}^{(3)} u_{2}=0 \tag{31}
\end{equation*}
$$

reduced separation occurs. Conditions (31) mean that the following separated equation for $S_{2}(\theta)$ holds

$$
\begin{equation*}
\frac{d^{2}}{d \theta^{2}} S_{2}+\left(\frac{d}{d \theta} S_{2}\right)^{2}=c_{2}, \quad c_{2} \in \mathbb{R} \tag{32}
\end{equation*}
$$

i.e., that $\psi_{2}(\theta)=e^{S_{2}}$ satisfies

$$
\begin{equation*}
\psi_{2}^{\prime \prime}=c_{2} \psi_{2} \tag{33}
\end{equation*}
$$

As in the Cartesian case plugging (32) into the separated equation $H_{s}$ provides a separated equation which involves the variable $r$ only:

$$
\begin{align*}
& r^{4} u_{1}^{(4)}+\left(4 r^{4} u_{1}+2 r^{3}\right) u_{1}^{(3)}+3 r^{4}\left(u_{1}^{(2)}\right)^{2}+r^{2}\left(6 r^{2}\left(u_{1}\right)^{2}+6 r u_{1}+2-1\right) u_{1}^{(2)}+ \\
& +r^{4}\left(u_{1}\right)^{4}+2 r^{3}\left(u_{1}\right)^{3}+\left(2 c_{2}-1\right)\left(r^{2}\left(u_{1}\right)^{2}-r u_{1}\right)-\lambda r^{4}+c_{2}^{2}-4 c_{2}=0 \tag{34}
\end{align*}
$$

here $c_{2}=u_{2}^{(2)}+\left(u_{2}\right)^{2}$ which is constant on the constraint surface $N_{p o l}$. The above equation becomes simpler passing to $\psi_{1}=e^{S_{1}}$

$$
\begin{equation*}
r^{4}\left(\psi_{1}\right)^{(4)}+2 r^{3}\left(\psi_{1}\right)^{(3)}+r\left(2 c_{2}-1\right)\left(r \psi_{1}^{\prime \prime}-\psi_{1}^{\prime}\right)+\left(c_{2}^{2}+4 c_{2}-\lambda r^{4}\right) \psi_{1}=0 \tag{35}
\end{equation*}
$$

whose solutions will be analyzed in Section 5.
The condition for the $r$-depending part of (30) is equivalent to

$$
r^{2}\left(\frac{d^{2}}{d r^{2}} S_{1}+\left(\frac{d}{d r} S_{1}\right)^{2}\right)-r \frac{d}{d r} S_{1}=c_{1}, \quad c_{1} \in \mathbb{R}
$$

However, this condition and its differential conditions defines a submanifold which is tangent to the generator $D_{1}$ only for $\lambda=0$.

### 4.3 Parabolic and Elliptic-Hyperbolic coordinates

We consider the general case of Liouville coordinates on $\mathbb{R}^{2}$ which produces Cartesian, polar, parabolic and elliptic-hyperbolic coordinates as special cases when we demand that the Gaussian curvature vanishes. This coordinate system is characterized by the following metric

$$
\begin{equation*}
g=(f(u)+g(v))\left(d u^{2}+d v^{2}\right) \tag{36}
\end{equation*}
$$

for some arbitrary smooth functions $f$ and $g$. We first consider the Helmholtz equation

$$
\begin{equation*}
\Delta \psi=\gamma \psi \tag{37}
\end{equation*}
$$

where we recall that $\Delta$ denotes the Laplace-Beltrami operator (8). In Liouville coordinates it takes the form

$$
\begin{equation*}
\Delta \psi=\frac{1}{f+g}\left(\partial_{u}^{2} \psi+\partial_{v}^{2} \psi\right) \tag{38}
\end{equation*}
$$

If we assume $\psi(u, v)=U(u) V(v)$ the Helmholtz equation may be written as

$$
\begin{equation*}
\left(\frac{U^{\prime \prime}}{U}-\gamma f\right)+\left(\frac{V^{\prime \prime}}{V}-\gamma g\right)=0 \tag{39}
\end{equation*}
$$

We observe that the equation separates for any smooth functions $f$ and $g$.

$$
\begin{equation*}
U^{\prime \prime}-(c+\gamma f) U=0 \quad V^{\prime \prime}+(c-\gamma g) V=0 \tag{40}
\end{equation*}
$$

where $c$ denotes the separation constant.
We now turn our attention to the bi-Helmoltz equation. The Gaussian curvature in this coordinate system has the following form

$$
\begin{equation*}
K=-\frac{1}{2(f+g)^{2}}\left(f^{\prime \prime}+g^{\prime \prime}-\frac{f^{\prime 2}+g^{\prime 2}}{f+g}\right) \tag{41}
\end{equation*}
$$

From here on we will set $K=0$ in order to consider the parabolic and elliptichyperbolic cases in which we also have that $f^{\prime} \neq 0, g^{\prime} \neq 0$. We use the LaplaceBeltrami (38) to write the bi-Laplace operator as

$$
\begin{equation*}
\Delta^{2} \psi=\frac{\Delta_{0}^{2} \psi}{(f+g)^{2}}-\frac{2}{(f+g)^{3}}\left(f^{\prime} \partial_{u} \Delta_{0} \psi+g^{\prime} \partial_{v} \Delta_{0} \psi\right)+\frac{f^{\prime \prime}+g^{\prime \prime}}{(f+g)^{3}} \Delta_{0} \psi \tag{42}
\end{equation*}
$$

From here we can write the bi-Helmholtz equation $\Delta^{2} \psi=\lambda \psi$, using (42). The product ansatz $\psi(u, v)=U(u) V(v)$ implies that (7) takes the form

$$
\begin{align*}
& f \frac{U^{(4)}}{U}+g \frac{V^{(4)}}{V}-2 f^{\prime} \frac{U^{(3)}}{U}-2 g^{\prime} \frac{V^{(3)}}{V}+f^{\prime \prime} \frac{U^{\prime \prime}}{U}+g^{\prime \prime} \frac{V^{\prime \prime}}{V} \\
+ & \frac{V^{\prime \prime}}{V}\left(f^{\prime \prime}+2 f \frac{U^{\prime \prime}}{U}-2 f^{\prime} \frac{U^{\prime}}{U}\right)+\frac{U^{\prime \prime}}{U}\left(g^{\prime \prime}+2 g \frac{V^{\prime \prime}}{V}-2 g^{\prime} \frac{V^{\prime}}{V}\right) \\
- & \lambda(f+g)^{3}=0 \tag{43}
\end{align*}
$$

We observe that the above equation is not separable and hence does not admit regular multiplicative separation. This result is consistent with the conclusion of Theorem 1. However, non-regular separation may be possible. By taking derivatives of (43) we can obtain necessary conditions for separation (constraint equations). We find a first necessary condition by applying $\partial_{u, v}^{2}$ which yields

$$
\begin{align*}
& f^{\prime}\left(\frac{V^{(4)}}{V}\right)^{\prime}+g^{\prime}\left(\frac{U^{(4)}}{U}\right)^{\prime}+\left(\frac{V^{\prime \prime}}{V}\right)^{\prime}\left(f^{\prime \prime}+2 f \frac{U^{\prime \prime}}{U}-2 f^{\prime} \frac{U^{\prime}}{U}\right)^{\prime} \\
& \quad+\left(\frac{U^{\prime \prime}}{U}\right)^{\prime}\left(g^{\prime \prime}+2 g \frac{V^{\prime \prime}}{V}-2 g^{\prime} \frac{V^{\prime}}{V}\right)^{\prime}-6 \lambda f^{\prime} g^{\prime}(f+g)=0 \tag{44}
\end{align*}
$$

We again note that this condition is not separable. We proceed by dividing by $f^{\prime} g^{\prime}$ and applying $\partial_{u, v}^{2}$ to get the separable condition

$$
\begin{align*}
& \left(\frac{1}{g^{\prime}}\left(\frac{V^{\prime \prime}}{V}\right)^{\prime}\right)^{\prime}\left(\frac{\left(f^{\prime \prime}+2 f \frac{U^{\prime \prime}}{U}-2 f^{\prime} \frac{U^{\prime}}{U}\right)^{\prime}}{f^{\prime}}\right)^{\prime} \\
& +\left(\frac{1}{f^{\prime}}\left(\frac{U^{\prime \prime}}{U}\right)^{\prime}\right)^{\prime}\left(\frac{\left(g^{\prime \prime}+2 g \frac{V^{\prime \prime}}{V}-2 g^{\prime} \frac{V^{\prime}}{V}\right)^{\prime}}{g^{\prime}}\right)^{\prime}=0 \tag{45}
\end{align*}
$$

From here we obtain several cases: (i) neither $\left(\frac{1}{g^{\prime}}\left(\frac{V^{\prime \prime}}{V}\right)^{\prime}\right)^{\prime}$ nor $\left(\frac{1}{g^{\prime}}\left(\frac{V^{\prime \prime}}{V}\right)^{\prime}\right)^{\prime}$ vanish, from which we obtain $\lambda=0$ (ii) only one vanishes which also implies $\lambda=0$ (iii) both vanish which gives us a solution to the Helmholtz equation in Liouville coordinates. Hence, we do not get any non-trivial separated solutions. The proof in case (i) is given below and in Appendix A while that for cases (ii) and (iii) is given in Appendix B.

Case(i): Dividing (45) by $\left(\frac{1}{f^{\prime}}\left(\frac{U^{\prime \prime}}{U}\right)^{\prime}\right)^{\prime}\left(\frac{1}{g^{\prime}}\left(\frac{V^{\prime \prime}}{V}\right)^{\prime}\right)^{\prime}$ we obtain a separable equation. Separating and integrating (see appendix A for this calculation) we find the following equations

$$
\begin{array}{r}
(2 g-C) \frac{V^{\prime \prime}}{V}-2 g^{\prime} \frac{V^{\prime}}{V}+g^{\prime \prime}-C_{1} g-C_{2}=0 \\
(2 f+C) \frac{U^{\prime \prime}}{U}-2 f^{\prime} \frac{U^{\prime}}{U}+f^{\prime \prime}-D_{1} f-D_{2}=0 \tag{47}
\end{array}
$$

where $C, C_{1}, C_{2}, D_{1} . D_{2} \in \mathbb{R}$. We can use these conditions in (44) to obtain

$$
\begin{equation*}
\frac{1}{g^{\prime}}\left(\frac{V^{(4)}}{V}\right)^{\prime}+\frac{1}{f^{\prime}}\left(\frac{U^{(4)}}{U}\right)^{\prime}+\frac{D_{1}}{g^{\prime}}\left(\frac{V^{\prime \prime}}{V}\right)^{\prime}+\frac{C_{1}}{f^{\prime}}\left(\frac{U^{\prime \prime}}{U}\right)^{\prime}-6 \lambda(f+g)=0 \tag{48}
\end{equation*}
$$

Notice that this condition is now separable. And using this condition we can integrate back to get a simplified form of equation (43)

$$
\begin{align*}
& f \frac{U^{(4)}}{U}+g \frac{V^{(4)}}{V}-2 f^{\prime} \frac{U^{(3)}}{U}-2 g^{\prime} \frac{V^{(3)}}{V}+f^{\prime \prime} \frac{U^{\prime \prime}}{U}+g^{\prime \prime} \frac{V^{\prime \prime}}{V}-\lambda f^{3}-\lambda g^{3}+ \\
& +f \frac{V^{(4)}}{V}+g \frac{U^{(4)}}{U}+\left(C_{1} g+C_{2}\right) \frac{U^{\prime \prime}}{U}+\left(D_{1} f+D_{2}\right) \frac{V^{\prime \prime}}{V}-3 \lambda f^{2} g-3 \lambda f g^{2}=0 \tag{49}
\end{align*}
$$

To make sure we account for all the constraints properly we also need to use the derivatives of equations (46) and (47). There are other constraints coming from the additional requirement $K=0$, for the interest of brevity we will include these calculations Appendix A.

After using all the constraint equations we end up with the following

$$
\begin{align*}
& \frac{\left(C_{2}-D_{2}-\left(C_{1}+D_{1}\right) f\right)}{C_{1}+D_{1}}\left(\alpha f+\beta+3 \lambda f^{2}-\frac{\left(k+C_{1}\right)\left(D_{1}-k\right) f+C_{1}\left(D_{2}-D\right)+D\left(D_{1}-k\right)}{2 f+C}\right. \\
& \left.+2 \frac{f^{\prime 2}}{(2 f+C)^{2}}\left(D_{1}-k\right)-\frac{\left(D_{1} f+D_{2}\right)^{2}-f^{\prime \prime 2}}{(2 f+C)^{2}}\right)-2 \frac{\left(D_{1}-k\right)}{2 f+C} f^{\prime 2} \\
& +\left(\left(k-C_{1}\right) f+C_{2}+D\right)\left(\frac{\left(D_{1}-k\right) f+D_{2}-D}{2 f+C}\right)+2 \lambda f^{3}+\alpha f^{2}+(\beta+\gamma) f=\delta \tag{50}
\end{align*}
$$

Where from the condition $K=0$ we have $f^{\prime 2}=k f^{2}+2 D f-k$ and $g^{\prime 2}=$ $-k g^{2}+2 D g+k$. Thus the above equation simplifies to a polynomial in $f$ after multiplication by $(2 f+C)^{2}$. Furthermore, since $f^{\prime} \neq 0$, the set $\left\{1, f, f^{2}, \ldots, f^{n}\right\}$ is linearly independent. This implies that the coefficients of the different powers of $f$ must all vanish. The coefficient of the highest power of $f$ the is $8 \lambda$ the vanishing of which implies that $\lambda=0$. In the case that $C_{1}+D_{1}=0$ (89) is a polynomial in $f$ with highest order term $3 \lambda f^{2}$ which also implies $\lambda=0$.

## 5 The circular vibrating plate

As mentioned in Section 3 one can pass between multiplicative and additive separation by the change of variable $\psi=e^{u}$. In this section we use the framework of multiplicative separation.

We apply the results for polar coordinates obtained in Section 4 to study vibrations of a thin, solid plate of uniform material, constant thickness, and areal density $\rho$ [11, 5].

The equation of motion for the plate under the assumption of small oscillations is given by [5]

$$
\begin{equation*}
\rho \ddot{\psi}+c \Delta^{2} \psi=0 \tag{51}
\end{equation*}
$$

where $\psi$ is the lateral elevation at any point on the region $\Omega$ in the interior of the plate and $c$ is some positive constant depending on the material. If the plate is clamped at the edge he boundary conditions on $\psi$ are $\left.\psi\right|_{\partial \Omega}=\left.\psi_{n}\right|_{\partial \Omega}=0$, where the subscript $n$ denotes differentiation with respect to the normal to the boundary. For the case of a circular plate $\Omega$ is a disk of some radius $a$. We obtain the solution of (51) by utilizing a separation of variables approach rather that the Fourier series method employed by Rayleigh [11].

Consider solutions of (51) of the form $\psi=w(r, \theta) T(t)$ which yield the separated equations

$$
\begin{equation*}
\frac{\Delta^{2} w}{w}=-\frac{\rho}{c} \frac{\ddot{T}}{T}=k^{4} . \tag{52}
\end{equation*}
$$

The separation constant is assumed to be positive (hence written as $k^{4}$ for convenience), since we expect oscillatory behaviour in time. Notice that this separation is regular since there are no additional constraints on the separated functions. The time part of (52) has the form

$$
\begin{equation*}
\ddot{T}+\omega^{2} T=0 \tag{53}
\end{equation*}
$$

where $\omega^{2}=\frac{c k^{4}}{\rho}$. The general solution is given by

$$
\begin{equation*}
T(t)=G \cos (\omega t)+H \sin (\omega t) \tag{54}
\end{equation*}
$$

where $G$ and $H$ are arbitrary constants.
The spatial part (52) satisfies the bi-Helmholtz equation,

$$
\begin{equation*}
\Delta^{2} w=k^{4} w \tag{55}
\end{equation*}
$$

Written out explicitly in polar coordinates this equation takes the form

$$
\begin{equation*}
w_{r r r r}+\frac{1}{r^{4}} w_{\theta \theta \theta \theta}+\frac{2}{r} w_{r r r}-\frac{1}{r^{2}} w_{r r}+\frac{1}{r^{3}} w_{r}+\frac{4}{r^{4}} w_{\theta \theta}-\frac{2}{r^{3}} w_{r \theta \theta}+\frac{2}{r^{2}} w_{r r \theta \theta}=k^{4} w, \tag{56}
\end{equation*}
$$

where the subscript $r$ denotes $\frac{d}{d r}$. We now consider product solutions of the form

$$
\begin{equation*}
w(r, \theta)=R(r) \Theta(\theta) \tag{57}
\end{equation*}
$$

which yields the separated equations

$$
\begin{gather*}
\Theta^{\prime \prime}=-l^{2} \Theta  \tag{58}\\
D^{4} R+\frac{2}{r} D^{3} R-\frac{1}{r^{2}}\left(1+2 l^{2}\right) D^{2} R+\frac{1}{r^{3}}\left(1+2 l^{2}\right) D R-\frac{l^{2}}{r^{4}}\left(4-l^{2}\right) R=k^{4} R \tag{59}
\end{gather*}
$$

where $l$ is a separation constant and $D=\frac{d}{d r}$.
The solution of the differential equation (58) satisfied by the angular function $\Theta$ has the form

$$
\begin{equation*}
\Theta_{l}(\theta)=E_{l} \cos (l \theta)+F_{l} \sin (l \theta) \tag{60}
\end{equation*}
$$

where $E_{l}$ and $F_{l}$ are arbitrary constants. In order that the function $\Theta$ be single valued on $\Omega$, the constant $l$ must be a positive integer

$$
l=n \quad(n=1,2, \ldots)
$$

We now turn our attention to (59). It may be shown that the left-hand-side may be factored in two ways to yield

$$
\begin{align*}
& \left(D^{2}+\frac{1}{r} D-\frac{n^{2}}{r^{2}}+k^{2}\right)\left(D^{2}+\frac{1}{r} D-\frac{n^{2}}{r^{2}}-k^{2}\right) R=0  \tag{61}\\
& \left(D^{2}+\frac{1}{r} D-\frac{n^{2}}{r^{2}}-k^{2}\right)\left(D^{2}+\frac{1}{r} D-\frac{n^{2}}{r^{2}}+k^{2}\right) R=0 \tag{62}
\end{align*}
$$

This result shows that the differential operators which appear inside the brackets (the Bessel and modified Bessel operators) commute. Since we know bases of the solution spaces to the Bessel and modified Bessel equations, the general solution of the radial equation may be written as

$$
\begin{equation*}
R_{n}(r)=A_{n} J_{n}(k r)+B_{n} Y_{n}(k r)+C_{n} I_{n}(k r)+D_{n} K_{n}(k r), \tag{63}
\end{equation*}
$$

where $J_{n}, Y_{n}, I_{n}, K_{n}$ are respectively the Bessel and modified Bessel functions of the first and second kind and $A_{n}, B_{n}, C_{n}, D_{n}$ are arbitrary constants.

The Bessel functions of the second kind are inadmissible since they are singular at the origin so we take $B_{n}=D_{n}=0$.

The boundary conditions tell us that $R_{n}(a)=R_{n}^{\prime}(a)=0$.

$$
\begin{align*}
& R_{n}(a)=A_{n} J_{n}(k a)+C_{n} I_{n}(k a)=0  \tag{64}\\
& R_{n}^{\prime}(a)=k A_{n} J_{n}^{\prime}(k a)+k C_{n} I_{n}^{\prime}(k a)=0 \tag{65}
\end{align*}
$$

To solve (64) we either need $k a=j_{n, m}$ (the $m$ th zero of $J_{n}$ ) and $C_{n}=0$ in which case we retrieve the vibrating membrane solution, or $A_{n}=-\frac{I_{n}(k a)}{J_{n}(k a)} C_{n}$ (we can absorb $C_{n}$ into the definition of $E_{n}, F_{n}$ to simplify). Since we are interested in solutions other than the ones for the vibrating membrane we take the latter in which case (65) becomes

$$
\begin{align*}
I_{n}^{\prime}(k a)-\frac{I_{n}(k a)}{J_{n}(k a)} J_{n}^{\prime}(k a) & =0  \tag{66}\\
\frac{I_{n}^{\prime}(k a)}{I_{n}(k a)}-\frac{J_{n}^{\prime}(k a)}{J_{n}(k a)} & =0 \tag{67}
\end{align*}
$$

We determine the roots of the above equation numerically to get a condition $k a=l_{n, m}$, where $l_{n, m}$ is the $m$ th root of equation (67) for some fixed $m$

$$
\begin{equation*}
R_{n, m}(r)=C_{n}\left(I_{n}\left(\frac{l_{n, m} r}{a}\right)-\frac{I_{n}\left(l_{n, m}\right)}{J_{n}\left(l_{n, m}\right)} J_{n}\left(\frac{l_{n, m} r}{a}\right)\right) \tag{68}
\end{equation*}
$$

Thus the general solution for $u$ has the form

$$
\begin{equation*}
\psi=\sum_{m, n}\left(E_{n} \cos (n \theta)+F_{n} \sin (n \theta)\right)\left(G_{n, m} \cos \left(\omega_{n, m} t\right)+H_{n, m} \sin \left(\omega_{n, m} t\right)\right) R_{n, m} \tag{69}
\end{equation*}
$$

Where $\omega_{n, m}=c^{2} \frac{l_{n, m}^{2}}{a^{2}}$. See Rayleigh [11] for further study of this solution.

## 6 Conclusion

The main idea of the paper is to apply the technique of the regular and nonregular separation to the search of the multiplicative separated solutions of the bi-Helmholtz equation, which is a classical 4 -th order PDE equation of Mathematical Physics which was solved in particular cases by means of this technique ([11]). Nevertheless, deeper studies about separability usually deal with first and second order PDE's (as well as the Hamilton-Jacobi equation and the stationary Schrödinger equation).

The choice of the bi-Helmholtz equation relies on two aspects: the existence of a physical application (it is not a simple toy model) and the fact that this kind of equation seemed to provide a good example for a deeper understanding of the non-regular separation method.

Non-regular separation appears naturally in the study of separability of the bi-Helmholtz equation, since we prove (Section 3) that regular separation never occurs for this equation, but on the other hand we already know the existence of a family of separated solutions (the trivial ones, arising from the separability of the standard Helmholtz equation). In the examples on the Euclidean plane, we show that, in two of the four separable coordinate systems (Cartesian and polar coordinates), non-trivial separated solutions can be determined, while in the remaining coordinate systems (parabolic and elliptic-hyperbolic coordinates) the only possibles separated solutions are proved to be the trivial ones.

The study of the separated solutions is done from the simplest geometrical view point: the search for a submanifold, as large as possible, where the separability conditions are satisfied, without assuming any particular structure for the separated solution (such as as side conditions or generalizations of Stäckel matrices as is done in [7]).

This paper represents the first step in a program to analyse interesting examples of non-regular separation As further extensions of the present paper, one could consider examples in higher dimensions or in different Riemannian manifolds, to get more physical applications. Indeed, a wider collection of examples is the natural starting point of a study to gain a deeper insight into the geometric conditions (for instance in terms of symmetry operators) of the non-regular
separation of the bi-Helmholtz equation, as has been done for the fixed energy R-separation for the Schrödinger equation [4, 7].

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## A

Dividing (44) by $\left(\frac{U^{(3)}}{f^{\prime} U}-\frac{U^{\prime \prime} U^{\prime}}{g^{\prime} U^{2}}\right)^{\prime}\left(\frac{V^{(3)}}{g^{\prime} V}-\frac{V^{\prime \prime} V^{\prime}}{g^{\prime} V^{2}}\right)^{\prime}$, we see that we can separate as follows

$$
\begin{equation*}
\frac{\left(\frac{\left(g^{\prime \prime}+2 g \frac{V^{\prime \prime}}{V}-2 g^{\prime} \frac{V^{\prime}}{V}\right)^{\prime}}{g^{\prime}}\right)^{\prime}}{\left(\frac{V^{(3)}}{g^{\prime} V}-\frac{V^{\prime \prime} V^{\prime}}{g^{\prime} V^{2}}\right)^{\prime}}=-\frac{\left(\frac{\left(f^{\prime \prime}+2 f \frac{U^{\prime \prime}}{U}-2 f^{\prime} \frac{U^{\prime}}{U}\right)^{\prime}}{f^{\prime}}\right)^{\prime}}{\left(\frac{U^{(3)}}{f^{\prime} U}-\frac{U^{\prime \prime} U^{\prime}}{f^{\prime} U^{2}}\right)^{\prime}}=C \tag{70}
\end{equation*}
$$

Where $C \in \mathbb{R}$. Separating and integrating

$$
\begin{align*}
\frac{\left(g^{\prime \prime}+2 g \frac{V^{\prime \prime}}{V}-2 g^{\prime} \frac{V^{\prime}}{V}\right)^{\prime}}{g^{\prime}} & =\frac{C}{g^{\prime}}\left(\frac{V^{\prime \prime}}{V}\right)^{\prime}+C_{1}  \tag{71}\\
\frac{\left(f^{\prime \prime}+2 f \frac{U^{\prime \prime}}{U}-2 f^{\prime} \frac{U^{\prime}}{U}\right)^{\prime}}{f^{\prime}} & =-\frac{C}{f^{\prime}}\left(\frac{U^{\prime \prime}}{U}\right)^{\prime}+D_{1} \tag{72}
\end{align*}
$$

The case where one of $\left(\frac{U^{(3)}}{f^{\prime} U}-\frac{U^{\prime \prime} U^{\prime}}{g^{\prime} U^{2}}\right)^{\prime},\left(\frac{V^{(3)}}{g^{\prime} V}-\frac{V^{\prime \prime} V^{\prime}}{g^{\prime} V^{2}}\right)^{\prime}$ vanish corresponds to setting $C=0$ in one of the above equations. Multiplying through by $f^{\prime}$ and $g^{\prime}$ respectively and integrating once again we obtain

$$
\begin{align*}
& g^{\prime \prime}+2 g \frac{V^{\prime \prime}}{V}-2 g^{\prime} \frac{V^{\prime}}{V}=C \frac{V^{\prime \prime}}{V}+C_{1} g+C_{2}  \tag{73}\\
& f^{\prime \prime}+2 f \frac{U^{\prime \prime}}{U}-2 f^{\prime} \frac{U^{\prime}}{U}=-C \frac{U^{\prime \prime}}{U}+D_{1} f+D_{2} \tag{74}
\end{align*}
$$

Where $C_{1}, C_{2}, D_{1} . D_{2} \in \mathbb{R}$. Or more compactly

$$
\begin{align*}
(2 g-C) \frac{V^{\prime \prime}}{V}-2 g^{\prime} \frac{V^{\prime}}{V}+g^{\prime \prime}-C_{1} g-C_{2} & =0  \tag{75}\\
(2 f+C) \frac{U^{\prime \prime}}{U}-2 f^{\prime} \frac{U^{\prime}}{U}+f^{\prime \prime}-D_{1} f-D_{2} & =0 \tag{76}
\end{align*}
$$

We can now separate (48)

$$
\begin{align*}
& \frac{U^{(4)}}{U}+C_{1} \frac{U^{\prime \prime}}{U}-3 \lambda f^{2}=\alpha f+\beta  \tag{77}\\
& \frac{V^{(4)}}{V}+D_{1} \frac{V^{\prime \prime}}{V}-3 \lambda g^{2}=-\alpha g+\gamma \tag{78}
\end{align*}
$$

The condition $K=0$ separates into

$$
\begin{equation*}
\frac{f^{(3)}}{f^{\prime}}=-\frac{g^{(3)}}{g^{\prime}}=k \tag{79}
\end{equation*}
$$

Where $k \in \mathbb{R}$ is a separation constant. Integrating once

$$
\begin{equation*}
f^{\prime \prime}=k f+D \quad g^{\prime \prime}=-k g+D \tag{80}
\end{equation*}
$$

These are extra conditions that have to be taken into account. We also need to account for the derivatives of equation (75) and (76)

$$
\begin{align*}
\frac{U^{\prime \prime}}{U} & =\frac{2 f^{\prime}}{2 f+C} \frac{U^{\prime}}{U}+\frac{D_{1} f+D_{2}-f^{\prime \prime}}{2 f+C}  \tag{81}\\
\frac{U^{(3)}}{U} & =\frac{D_{1} f+D_{2}+f^{\prime \prime}}{2 f+C} \frac{U^{\prime}}{U}+\frac{D_{1}-k}{2 f+C} f^{\prime}  \tag{82}\\
\frac{U^{(4)}}{U} & =\frac{D_{1}}{2 f+C} f^{\prime \prime}+\frac{D_{1} f+D_{2}+f^{\prime \prime}}{2 f+C} \frac{U^{\prime \prime}}{U}-\frac{2 f^{\prime}}{2 f+C} \frac{U^{(3)}}{U}  \tag{83}\\
& =\frac{\left(D_{1}-k\right) f^{\prime \prime}}{2 f+C}+\frac{\left(D_{1} f+D_{2}\right)^{2}-f^{\prime \prime 2}}{(2 f+C)^{2}}-\frac{2 f^{\prime 2}}{(2 f+C)^{2}}\left(D_{1}-k\right)  \tag{84}\\
\frac{V^{\prime \prime}}{V} & =\frac{2 g^{\prime}}{2 g-C} \frac{V^{\prime}}{V}+\frac{C_{1} g+C_{2}+g^{\prime \prime}}{2 g-C}  \tag{85}\\
\frac{V^{(3)}}{V} & =\frac{C_{1} g+C_{2}+g^{\prime \prime}}{2 g-C} \frac{V^{\prime}}{V}+\frac{C_{1}+k}{2 g-C} g^{\prime}  \tag{86}\\
\frac{V^{(4)}}{V} & =\frac{C_{1}}{2 g-C} g^{\prime \prime}+\frac{C_{1} g+C_{2}-g^{\prime \prime}}{2 g-C} \frac{V^{\prime \prime}}{V}-\frac{2 g^{\prime}}{2 g-C} \frac{V^{(3)}}{V}  \tag{87}\\
& =\frac{\left(C_{1}+k\right) g^{\prime \prime}}{2 g+C}+\frac{\left(C_{1} g+C_{2}\right)^{2}-g^{\prime \prime 2}}{(2 g-C)^{2}}-\frac{2 g^{\prime 2}}{(2 g-C)^{2}}\left(C_{1}+k\right) \tag{88}
\end{align*}
$$

Using (81) to eliminate higher derivatives in (77)

$$
\begin{align*}
& \left(C_{1}+D_{1}\right) \frac{2 f^{\prime}}{2 f+C} \frac{U^{\prime}}{U}-2 \frac{f^{\prime 2}}{(2 f+C)^{2}}\left(D_{1}-k\right)+\left(D_{1}-k\right) \frac{f^{\prime \prime}}{2 f+C}+\frac{\left(D_{1} f+D_{2}\right)^{2}-f^{\prime \prime 2}}{(2 f+C)^{2}}+ \\
+ & C_{1} \frac{D_{1} f+D_{2}-f^{\prime \prime}}{2 f+C}-3 \lambda f^{2}=\alpha f+\beta  \tag{89}\\
& \left(C_{1}+D_{1}\right) \frac{2 g^{\prime}}{2 g-C} \frac{V^{\prime}}{V}-2 \frac{g^{\prime 2}}{(2 g-C)^{2}}\left(C_{1}+k\right)+\left(C_{1}+k\right) \frac{g^{\prime \prime}}{2 g-C}+\frac{\left(C_{1} g+C_{2}\right)^{2}-g^{\prime \prime 2}}{(2 g-C)^{2}}+ \\
+ & D_{1} \frac{C_{1} g+C_{2}-g^{\prime \prime}}{2 g-C}-3 \lambda g^{2}=-\alpha g+\gamma \tag{90}
\end{align*}
$$

Substituting (77) into (49) we have the following

$$
\begin{align*}
& -2 f^{\prime} \frac{U^{(3)}}{U}-2 g^{\prime} \frac{V^{(3)}}{V}+\left(\left(k-C_{1}\right) f+C_{2}+D\right) \frac{U^{\prime \prime}}{U}+\left(D_{2}+D-\left(k+D_{1}\right) g\right) \frac{V^{\prime \prime}}{V} \\
& +2 \lambda\left(f^{3}+g^{3}\right)+\alpha\left(f^{2}-g^{2}\right)+(\beta+\gamma)(f+g)=0 \tag{91}
\end{align*}
$$

Separating this equation we have

$$
\begin{gather*}
-2 f^{\prime} \frac{U^{(3)}}{U}+\left(\left(k-C_{1}\right) f+C_{2}+D\right) \frac{U^{\prime \prime}}{U}+2 \lambda f^{3}+\alpha f^{2}+(\beta+\gamma) f=\delta  \tag{92}\\
-2 g^{\prime} \frac{V^{(3)}}{V}+\left(-\left(k+D_{1}\right) g+D_{2}+D\right) \frac{V^{\prime \prime}}{V}+2 \lambda g^{3}-\alpha g^{2}+(\beta+\gamma) g=-\delta \tag{93}
\end{gather*}
$$

For some $\delta \in \mathbb{R}$. Eliminating the derivatives from (92) we have the following

$$
\begin{align*}
& 2 f^{\prime} \frac{C_{2}-D_{2}-\left(C_{1}+D_{1}\right) f}{2 f+C} \frac{U^{\prime}}{U}-2 \frac{\left(D_{1}-k\right)}{2 f+C} f^{\prime 2}  \tag{94}\\
& +\left(\left(k-C_{1}\right) f+C_{2}+D\right)\left(\frac{\left(D_{1}-k\right) f+D_{2}-D}{2 f+C}\right) \\
& +2 \lambda f^{3}+\alpha f^{2}+(\beta+\gamma) f=\delta \tag{95}
\end{align*}
$$

Isolating for $\frac{2 f^{\prime}}{2 f+C} \frac{U^{\prime}}{U}$ from (89) assuming $C_{1}+D_{1} \neq 0$ and using this to eliminate derivatives in (94)

$$
\begin{align*}
& \frac{\left(C_{2}-D_{2}-\left(C_{1}+D_{1}\right) f\right)}{C_{1}+D_{1}}\left(\alpha f+\beta+3 \lambda f^{2}-\frac{\left(k+C_{1}\right)\left(D_{1}-k\right) f+C_{1}\left(D_{2}-D\right)+D\left(D_{1}-k\right)}{2 f+C}\right. \\
& \left.+2 \frac{f^{\prime 2}}{(2 f+C)^{2}}\left(D_{1}-k\right)-\frac{\left(D_{1} f+D_{2}\right)^{2}-f^{\prime \prime 2}}{(2 f+C)^{2}}\right)-2 \frac{\left(D_{1}-k\right)}{2 f+C} f^{\prime 2} \\
& +\left(\left(k-C_{1}\right) f+C_{2}+D\right)\left(\frac{\left(D_{1}-k\right) f+D_{2}-D}{2 f+C}\right)+2 \lambda f^{3}+\alpha f^{2}+(\beta+\gamma) f=\delta \tag{96}
\end{align*}
$$

When $C_{1}+D_{1}=0$, (89) implies $\lambda=0$.

## B

Case (ii): In this case we have

$$
\begin{equation*}
\left(\frac{1}{f^{\prime}}\left(\frac{U^{\prime \prime}}{U}\right)^{\prime}\right)^{\prime}=0, \quad\left(\frac{1}{g^{\prime}}\left(\frac{V^{\prime \prime}}{V}\right)^{\prime}\right)^{\prime} \neq 0 \tag{97}
\end{equation*}
$$

Integrating we get

$$
\begin{equation*}
\frac{U^{\prime \prime}}{U}=\alpha_{5} f+\alpha_{6} \tag{98}
\end{equation*}
$$

where $\alpha_{5}, \alpha_{6}$ are constants. Equation (45) implies that

$$
\begin{equation*}
2 f \frac{U^{\prime \prime}}{U}-2 f^{\prime} \frac{U^{\prime}}{U}+f^{\prime \prime}=\alpha_{3} f+\alpha_{4} \tag{99}
\end{equation*}
$$

$\alpha_{3}, \alpha_{4}$ are constants. Substitution for $\frac{U^{\prime \prime}}{U}$ from (98) yields

$$
\begin{equation*}
2 f^{\prime} U^{\prime}=\left(f^{\prime \prime}+2 \alpha_{5} f^{2}+\left(2 \alpha_{6}-\alpha_{3}\right) f-\alpha_{4}\right) U \tag{100}
\end{equation*}
$$

Differentiation of the above equation followed substitution for $U^{\prime \prime}$ from (98) and $2 f^{\prime} U^{\prime}$ from (100) yields after simplification

$$
\begin{equation*}
\left(2 \alpha_{5} f^{2}+\left(2 \alpha_{6}-\alpha_{3}\right) f-\alpha_{4}^{2}\right)^{2}-f^{\prime \prime 2}+2 f^{\prime} f^{(3)}+4 \alpha_{5} f f^{\prime 2}-\alpha_{3} f^{\prime 2}=0 \tag{101}
\end{equation*}
$$

Using the relations between $f$ and its derivatives

$$
\begin{equation*}
f^{\prime \prime}=k f+D, \quad f^{\prime 2}=k f^{2}+2 D f+\Lambda \tag{102}
\end{equation*}
$$

Equation (101) becomes a polynomial in $f$, the coefficient of the highest power of $f$ is $4 \alpha_{5}^{2}$ which implies that $\alpha_{5}=0$. In view of the above equation (45) separates, the compatibility of the separated equation for $U$ with (98) and (100) gives us that $\lambda=\alpha_{5}^{2}$, thus we conclude that $\lambda=0$ in this case as well.

Case (iii): both of the following conditions hold:

$$
\begin{align*}
& \left(\frac{1}{f^{\prime}}\left(\frac{U^{\prime \prime}}{U}\right)^{\prime}\right)^{\prime}=0  \tag{103}\\
& \left(\frac{1}{g^{\prime}}\left(\frac{V^{\prime \prime}}{V}\right)^{\prime}\right)^{\prime}=0 \tag{104}
\end{align*}
$$

The solutions of (103) and (104) are given by

$$
\begin{gather*}
\frac{U^{\prime \prime}}{U}=\alpha_{5} f+\alpha_{6}  \tag{105}\\
\frac{V^{\prime \prime}}{V}=\left(\beta_{5} f+\beta_{6}\right) \tag{106}
\end{gather*}
$$

Computing derivatives of the above equations:

$$
\begin{align*}
& U^{(3)}=\left(\alpha_{5} f+\alpha_{6}\right) U^{\prime}+\alpha_{5} f^{\prime} U  \tag{107}\\
& U^{(4)}=2 \alpha_{5} f^{\prime} U^{\prime}+\left(\alpha_{5} f^{\prime \prime}+\left(\alpha_{5} f+\alpha_{6}\right)^{2}\right) U  \tag{108}\\
& V^{(3)}=\left(\beta_{5} g+\beta_{6}\right) V^{\prime}+\beta_{5} g^{\prime} V  \tag{109}\\
& V^{(4)}=2 \beta_{5} g^{\prime} V^{\prime}+\left(\beta_{5} g^{\prime \prime}+\left(\beta_{5} g+\beta_{6}\right)^{2}\right) V \tag{110}
\end{align*}
$$

With the use of the above derivatives the integrability condition (44) separates to yield the following equations:

$$
\begin{align*}
& 2\left(\beta_{5}-\alpha_{5}\right) f^{\prime} U^{\prime}=\left(\left(\alpha_{5}+\beta_{5}\right) f^{\prime \prime}+\alpha_{5}\left(\alpha_{5}+2 \beta_{5}\right) f^{2}\right.  \tag{111}\\
& \left.+\left(2 \alpha_{6}\left(\alpha_{5}+\beta_{5}\right)-\alpha\right) f-3 \lambda f^{2}+\alpha_{6}^{2}-\alpha_{7}\right) U  \tag{112}\\
& 2\left(\alpha_{5}-\beta_{5}\right) g^{\prime} V^{\prime}=\left(\left(\alpha_{5}+\beta_{5}\right) g^{\prime \prime}+\beta_{5}\left(\alpha_{5}+2 \beta_{5}\right) g^{2}\right.  \tag{113}\\
& \left.+\left(2 \beta_{6}\left(\alpha_{5}+\beta_{5}\right)+\alpha\right) g-3 \lambda g^{2}+\alpha_{6}^{2}-\alpha_{7}\right) V \tag{114}
\end{align*}
$$

where $\alpha$ is the separation constant. If $\beta_{5}=\alpha_{5}$, (111) and (113) imply that

$$
\begin{equation*}
\lambda=\alpha_{5}^{2}, \quad \beta_{5}=\alpha_{5}, \quad \beta_{6}=-\alpha_{6} \tag{115}
\end{equation*}
$$

We conclude that (115) implies that $\phi(u, v)=U(u) V(v)$ defines a separable solution of the Helmholtz equation (See (40).). If $\beta_{5} \neq \alpha_{5}$, one differentiates (111) and (113) and eliminates all derivatives of $U$ and $V$. One obtains polynomial equations in $f$ and $g$ which imply that $\beta_{5}^{2}=\alpha_{5}^{2}$. The case $\beta_{5}=\alpha_{5}$ has already been considered. The case $\beta_{5}=-\alpha_{5}$, yields $3 \lambda=-\alpha_{5}^{2}$, which is un-physical. This completes the proof of Case (iii).

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