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# Uniform theory of inhomogeneous waveguide modes near cut-off 

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#### Abstract

A uniform asymptotic method is used to calculate the modal eigenvalues for an inhomogeneous dielectric waveguide for modes near cut-off, thus extending previous work on uniform methods. The existence of a finite homogeneous cladding medium is properly accounted for, with arbitrary variations of core refractive index.


## 1. Introduction

In the preceding paper (Arnold 1980) a study was carried out of the construction of asymptotic approximations to the eigenvalue of the differential equation

$$
\begin{equation*}
\mathrm{d}^{2} \phi / \mathrm{d} \rho^{2}+\left(U^{2}-V^{2} f+\mu / \rho^{2}\right) \phi=0 \tag{1.1}
\end{equation*}
$$

where

$$
\begin{equation*}
\mu=\frac{1}{4}-m^{2}, \quad m \in\{0,1,2, \ldots\} \tag{1.2}
\end{equation*}
$$

$f$ is an arbitrary analytic function of $p^{2}$ such that

$$
\begin{array}{ll}
f=0, & \rho=0, \\
f=1, & \rho=1, \tag{1.3b}
\end{array}
$$

$V$ is a large parameter and $U^{2}$ is the eigenvalue which is required to be found. This differential equation arises in the theory of wave propagation on a dielectric cylinder having a radial variation of refractive index. Boundary conditions

$$
\begin{equation*}
\phi \sim \rho^{m+1 / 2}, \quad \rho \rightarrow 0, \tag{i}
\end{equation*}
$$

(ii) $\quad \mathrm{d} \phi / \mathrm{d} \rho=K \phi, \quad \rho=1$
have to be applied, where

$$
\begin{align*}
& K=\frac{1}{2}+W K_{m}^{\prime}(W) / \mathrm{K}_{m}(W),  \tag{1.5}\\
& W^{2}=V^{2}-U^{2} \tag{1.6}
\end{align*}
$$

and $\mathrm{K}_{m}(W)$ and $\mathrm{K}_{m}^{\prime}(W)$ are the modified Hankel function and its derivative respectively. By assuming that $U^{2}$ was not too close to $V^{2}$ and $\mu$ was small compared to $V^{2}$, it was possible to obtain an expression for the eigenvalue $U^{2}$ (as $V \rightarrow \infty$ ) by constructing an asymptotic approximation to $\phi$ which is valid on all $0 \leqslant \rho \leqslant 1$. This is in contrast to
the approach of Kurtz and Streifer (1969), who used matched piecewise-uniform asymptotic representations. A perturbation theory for quasi-quadratic $f$,

$$
f=\rho^{2}(1+\epsilon g)
$$

was outlined, where $\epsilon$ is a small parameter and $g$ an arbitrary analytic function of $\rho^{2}$.
However, the case of near cut-off modes, when $U^{2} \sim V^{2}$, was not considered, and this case introduces complications in the determination of the eigenvalue which we wish to examine in detail here. In the JWKB analysis of wave propagation problems (Fröman and Fröman 1965) certain points, called turning points in quantum mechanics and caustics in geometrical optics, are known to play a crucial role; at such points JWKB theory fails to give a valid asymptotic representation. In this problem, the caustic is the zero of the leading order term in the coefficient of $\phi$ in equation (1.1) (as $V \rightarrow \infty$ ), that is to say a point $\rho_{1}$ at which

$$
\begin{equation*}
U^{2}=V^{2} f, \quad \rho=\rho_{1} \tag{1.7}
\end{equation*}
$$

Since $f=1$ at $\rho=1$, it is apparent from (1.7) that, as $U^{2} \rightarrow V^{2}, \rho_{1} \rightarrow 1$. Thus there is a confluence of the caustic $\rho=\rho_{1}$ and the end-point $\rho=1$ at which the boundary condition (1.4b) has to be applied. This invalidates the procedure used previously (Arnold 1980), as it was there assumed that the caustic and the boundary point were sufficiently well separated to allow the use of the Liouville-Green approximation (Olver 1974) to approximate $\phi$ at the boundary, this approximation subsequently being obtained from the original uniform representation. This means that higher-order approximation functions (Airy functions) are required in order to approximate the eigenvalue, and we shall seek to obtain these also from a uniform representation.

To do this it is convenient to make some small changes to the original differential equation (1.1). For this purpose we use (1.6) to obtain

$$
\begin{equation*}
\mathrm{d}^{2} \phi / \mathrm{d} \rho^{2}+\left(V^{2} Q^{2}-W^{2}+\mu / \rho^{2}\right) \phi=0 \tag{1.8}
\end{equation*}
$$

where

$$
\begin{equation*}
Q^{2}=1-f \tag{1.9}
\end{equation*}
$$

We now regard $W^{2}$ as an eigenvalue of (1.8) and assert that $Q^{2}$ satisfies the following conditions:

$$
\begin{equation*}
Q^{2}=1 \text { when } \rho=0 \tag{i}
\end{equation*}
$$

(ii) $\quad Q^{2}=0$ when $\rho=1$;
(iii) $\quad Q^{2}$ is an analytic function of $\rho^{2}$;
(iv) $\quad V^{2} Q^{2}-W^{2}$ has one zero in $0 \leqslant \rho \leqslant 1$;
(v) $\quad V^{2} Q^{2}-W^{2}+\mu / \rho^{2}$ has two zeros in $0 \leqslant \rho \leqslant 1$.

Furthermore, it is convenient to express the fact that $U^{2}$ is nearly equal to $V^{2}$ by allowing $W^{2}$ to be small;

$$
\begin{equation*}
\lambda=W^{2} / V^{2} \sim \mathrm{O}\left(V^{-\kappa}\right), \quad \kappa>0 \tag{1.11}
\end{equation*}
$$

is introduced as a hypothesis. This then ensures that the terms $V^{2} Q^{2}$ and $W^{2}$ have different asymptotic orders in $V$. The exponent $\kappa$ in general depends on the boundary conditions, but in this case it transpires that $\kappa=\frac{2}{3}$ is correct; we shall simply assume that

$$
\lambda \sim \mathrm{O}\left(V^{-2 / 3}\right)
$$

## 2. Uniform asymptotic solution

In the preceding paper we showed how a Liouville transform

$$
\begin{align*}
\left(z_{1}^{2}-z^{2}\right)(\mathrm{d} z / \mathrm{d} \rho)^{2} & =U^{2} / V^{2}-f  \tag{2.1a}\\
& =Q^{2}-W^{2} / V^{2} \tag{2.1b}
\end{align*}
$$

could be constructed, such that

$$
\begin{align*}
\int_{0}^{z_{1}}\left(z_{1}^{2}-z^{2}\right)^{1 / 2} \mathrm{~d} z & =\int_{0}^{\rho_{1}}\left(\frac{U^{2}}{V^{2}}-f\right)^{1 / 2} \mathrm{~d} \rho  \tag{2.2a}\\
& =\int_{0}^{\rho_{1}}\left(Q^{2}-\frac{W^{2}}{V^{2}}\right)^{1 / 2} \mathrm{~d} \rho \tag{2.2b}
\end{align*}
$$

where $\rho_{1}$ is the zero of the integrand. The point $z$ is given in terms of $\rho$ by

$$
\begin{equation*}
\int_{z}^{z_{1}}\left(z_{1}^{2}-z^{\prime 2}\right)^{1 / 2} \mathrm{~d} z^{\prime}=\int_{\rho}^{\rho_{1}}\left(Q^{2}-\lambda\right)^{1 / 2} \mathrm{~d} \rho^{\prime}, \tag{2.3}
\end{equation*}
$$

where

$$
\begin{equation*}
\lambda=W^{2} / V^{2} . \tag{2.4}
\end{equation*}
$$

In this sequel we will construct this transform in an indirect way which has the advantage of giving an explicit expansion equivalent to (2.3) (which is implicit); evaluation of this expansion is quite simple when $\rho=1$, which we need to apply the boundary conditions.

We begin by introducing the intermediate variable $\zeta$ :

$$
\begin{equation*}
\left(\zeta_{0}^{2}-\zeta^{2}\right)(\mathrm{d} \zeta / \mathrm{d} \rho)^{2}=Q^{2} \tag{2.5}
\end{equation*}
$$

In order that $\zeta$ and $\rho$ be analytic functions of each other it is necessary to ensure that the zeros of both sides of equation (2.5) are consistent; therefore

$$
\begin{equation*}
\int_{0}^{\zeta_{0}}\left(\zeta_{0}^{2}-\zeta^{2}\right)^{1 / 2} \mathrm{~d} \zeta=\int_{0}^{1} Q \mathrm{~d} \rho \tag{2.6}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{\pi}{4} \zeta_{0}^{2}=\int_{0}^{1} Q \mathrm{~d} \rho \tag{2.7}
\end{equation*}
$$

Next we suppose that

$$
\begin{equation*}
z=\zeta+\lambda \eta_{1}+\lambda^{2} \eta_{2}+\mathbf{O}\left(\lambda^{3}\right) . \tag{2.8a}
\end{equation*}
$$

This form for $z$ suggests itself as an extension of the method used by Olver (1975) for the case $\mu=0$.

If $z$ is to be the Liouville transform (through $\zeta$ ) of $\rho$, then $\eta_{1}$ and $\eta_{2}$ have to be chosen so that

$$
\begin{equation*}
\left(z_{1}^{2}-z^{2}\right)(\mathrm{d} z / \mathrm{d} \rho)^{2}=Q^{2}-\lambda, \tag{2.8b}
\end{equation*}
$$

and therefore

$$
\begin{align*}
\left(z_{1}^{2}-z^{2}\right)(\mathrm{d} z / \mathrm{d} \zeta)^{2} & =\left(Q^{2}-\lambda\right)(\mathrm{d} \rho / \mathrm{d} \zeta)^{2}  \tag{2.9}\\
& =\zeta_{0}^{2}-\zeta^{2}-\lambda(\mathrm{d} \rho / \mathrm{d} \zeta)^{2} \tag{2.10}
\end{align*}
$$

where (2.5) has been used. If we suppose further that

$$
\begin{equation*}
z_{1}^{2}=\alpha_{0}+\lambda \alpha_{1}+\lambda^{2} \alpha_{2}+\mathbf{O}\left(\lambda^{3}\right) \tag{2.11}
\end{equation*}
$$

then we may substitute (2.8b) and (2.11) in (2.10), and equate powers of $\lambda$ to find

$$
\begin{align*}
& \alpha_{0}=\zeta_{0}^{2}=\frac{4}{\pi} \int_{0}^{1} Q \mathrm{~d} \rho  \tag{2.12a}\\
& \zeta \eta_{1}-\left(\zeta_{0}^{2}-\zeta^{2}\right) \mathrm{d} \eta_{1} / \mathrm{d} \zeta=\frac{1}{2}\left[(\mathrm{~d} \rho / \mathrm{d} \zeta)^{2}+\alpha_{1}\right] \tag{2.12b}
\end{align*}
$$

and so on for higher-order terms. Equation (2.12b) is a differential equation for $\eta_{1}$ which is easily solved:

$$
\begin{equation*}
\eta_{1}=\frac{1}{2\left(\zeta_{0}^{2}-\zeta^{2}\right)^{1 / 2}} \int_{\zeta}^{\zeta_{0}}\left[\left(\frac{\mathrm{~d} \rho}{\mathrm{~d} \zeta^{\prime}}\right)^{2}+\alpha_{1}\right] \frac{\mathrm{d} \zeta^{\prime}}{\left(\zeta_{0}^{2}-\zeta^{\prime 2}\right)^{1 / 2}} \tag{2.13}
\end{equation*}
$$

The parameter $\alpha_{1}$, as yet undetermined, may be chosen to make $\eta_{1}$ vanish at the origin. In fact, this ensures that $\eta_{1}$ has only odd powers of $\zeta$ in its expansion about $\zeta=0$, which is necessary in order that $z$ and $\zeta$ be analytic functions of each other at $\zeta=0$ for any $\lambda$. Thus

$$
\begin{equation*}
\alpha_{1} \int_{0}^{\zeta_{0}} \frac{\mathrm{~d} \zeta^{\prime}}{\left(\zeta_{0}^{2}-\zeta^{\prime 2}\right)^{1 / 2}}=-\int_{0}^{\zeta_{0}}\left(\frac{\mathrm{~d} \rho}{\mathrm{~d} \zeta^{\prime}}\right)^{2} \frac{\mathrm{~d} \zeta^{\prime}}{\left(\zeta_{0}^{2}-\zeta^{\prime 2}\right)^{1 / 2}} \tag{2.14}
\end{equation*}
$$

and using (2.5) on the right of (2.14) gives

$$
\begin{equation*}
\frac{\pi}{2} \alpha_{1}=-\int_{0}^{1} \frac{\mathrm{~d} \rho}{Q} \tag{2.15}
\end{equation*}
$$

Similar considerations hold for $\alpha_{2}$ and $\eta_{2}$, but, as we shall not use them, we need not dwell on them; the differential equation for $\eta_{2}$ is

$$
\begin{equation*}
\zeta \eta_{2}-\left(\zeta_{0}^{2}-\zeta^{2}\right) \frac{\mathrm{d} \eta_{2}}{\mathrm{~d} \zeta}=\left(\zeta_{0}^{2}-\zeta^{2}\right)\left(\frac{\mathrm{d} \eta_{1}}{\mathrm{~d} \zeta}\right)^{2}-\eta_{1}^{2}+\alpha_{2}+\frac{2 \mathrm{~d} \eta_{1}}{\mathrm{~d} \zeta}\left(\alpha_{1}-2 \zeta \eta_{1}\right) \tag{2.16}
\end{equation*}
$$

and $\alpha_{2}$ is given by

$$
\begin{equation*}
\frac{\pi}{2} \alpha_{2}=\rho_{1}^{\prime} \alpha_{1}-\frac{1}{4} \int_{0}^{1}\left(1-\rho_{1}^{\prime} \frac{\mathrm{d} Q^{2}}{\mathrm{~d} \rho}\right) \frac{\mathrm{d} \rho}{Q^{3}} \tag{2.17}
\end{equation*}
$$

where

$$
\begin{equation*}
\rho_{1}^{\prime}=\lim _{\lambda \rightarrow 0}\left(\frac{\mathrm{~d} \rho_{1}}{\mathrm{~d} \lambda}\right) . \tag{2.18}
\end{equation*}
$$

(Equation (2.17) is obtained by expanding (2.2b) in powers of $\lambda$. The integral exists because, as $\rho \rightarrow 1$, the term in brackets vanishes.) In principle, all higher-order terms subsumed under the $\mathrm{O}\left(\lambda^{3}\right)$ terms could be obtained in this way. Instead, we make a small modification; we drop the term $\mathrm{O}\left(\lambda^{3}\right)$ altogether from ( $2.8 b$ ), so that we have exactly

$$
\begin{equation*}
z=\zeta+\lambda \eta_{1}+\lambda^{2} \eta_{2} \tag{2.19}
\end{equation*}
$$

Then equation (2.10) will need to be modified to

$$
\begin{equation*}
\left(z_{1}^{2}-z^{2}\right)(\mathrm{d} z / \mathrm{d} \zeta)^{2}=\zeta_{0}^{2}-\zeta^{2}-\lambda(\mathrm{d} \rho / \mathrm{d} \zeta)^{2}+\mathrm{O}\left(\lambda^{3}\right) \tag{2.20a}
\end{equation*}
$$

and ( $2.8 b$ ) becomes

$$
\begin{equation*}
\left(z_{1}^{2}-z^{2}\right)(\mathrm{d} z / \mathrm{d} \rho)^{2}=Q^{2}-\lambda+\mathrm{O}\left(\lambda^{3}\right) \tag{2.20b}
\end{equation*}
$$

Since $\rho, \zeta$ and $z$ are all analytic functions of each other, the $\mathrm{O}\left(\lambda^{3}\right)$ terms are uniform on $0 \leqslant \rho \leqslant 1$, and may be transferred to the left-hand side of (2.20b) to give

$$
\begin{equation*}
\left[z_{1}^{2}-z^{2}+\mathrm{O}\left(\lambda^{3}\right)\right](\mathrm{d} z / \mathrm{d} \rho)^{2}=Q^{2}-\lambda \tag{2.21}
\end{equation*}
$$

As previously, $z_{1}$ is given by (2.11). We may now absorb the $\mathrm{O}\left(\lambda^{3}\right)$ term in (2.11) into the corresponding term in (2.21), and define $z_{1}$ exactly by

$$
\begin{equation*}
z_{1}^{2}=\alpha_{0}+\lambda \alpha_{1}+\lambda^{2} \alpha_{2} . \tag{2.22}
\end{equation*}
$$

This completes our analysis of the transformation of the independent variable.
Now we transform $\phi$ by

$$
\begin{equation*}
\phi=(\mathrm{d} z / \mathrm{d} \rho)^{-1 / 2} \Phi . \tag{2.23}
\end{equation*}
$$

Then $\Phi$ satisfies

$$
\begin{equation*}
\mathrm{d}^{2} \Phi / \mathrm{d} z^{2}+\left\{V^{2}\left[z_{1}^{2}-z^{2}+\mathrm{O}\left(\lambda^{3}\right)\right]+\left(\mu / \rho^{2}\right)(\mathrm{d} \rho / \mathrm{d} z)^{2}+h\right\} \Phi=0, \tag{2.24}
\end{equation*}
$$

where

$$
\begin{equation*}
h=-\left(\frac{\mathrm{d} \rho}{\mathrm{~d} z}\right)^{1 / 2} \frac{\mathrm{~d}^{2}}{\mathrm{~d} z^{2}}\left(\frac{\mathrm{~d} \rho}{\mathrm{~d} z}\right)^{-1 / 2} . \tag{2.25}
\end{equation*}
$$

As in the preceding paper it can be shown that

$$
\begin{equation*}
\frac{1}{\rho^{2}}\left(\frac{\mathrm{~d} \rho}{\mathrm{~d} z}\right)^{2}=\frac{1}{z^{2}}+\mathrm{O}(1) \tag{2.26}
\end{equation*}
$$

and so, if $\mu \sim \mathrm{O}(1)$,

$$
\begin{equation*}
\frac{\mu}{\rho^{2}}\left(\frac{\mathrm{~d} \rho}{\mathrm{~d} z}\right)^{2}=\frac{\mu}{z^{2}}+\mathrm{O}(1) \tag{2.27}
\end{equation*}
$$

The symbol $\mathrm{O}(1)$ refers to a function of $z$ which is analytic on $0 \leqslant \rho \leqslant 1$ and is $\mathrm{O}(1)$ as $V \rightarrow \infty$.

Similarly, $h$ is an $\mathrm{O}(1)$ analytic function, and $V^{2} \mathrm{O}\left(\lambda^{3}\right) \sim \mathrm{O}(1)$ if $\lambda \sim \mathrm{O}\left(V^{-2 / 3}\right)$. All these terms may be collected together to give

$$
\begin{equation*}
\mathrm{d}^{2} \Phi / \mathrm{d} z^{2}+\left[V^{2}\left(z_{1}^{2}-z^{2}\right)+\mu / z^{2}+\delta\right] \Phi=0, \tag{2.28}
\end{equation*}
$$

where

$$
\begin{equation*}
\delta \sim \mathrm{O}(1) \tag{2.29}
\end{equation*}
$$

As before (Arnold 1980) we can show that, as $V \rightarrow \infty$,

$$
\begin{equation*}
\Phi \sim \Phi_{0} \tag{2.30}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathrm{d}^{2} \Phi_{0} / \mathrm{d} z^{2}+\left[V^{2}\left(z_{1}^{2}-z^{2}\right)+\mu / z^{2}\right] \Phi_{0}=0 \tag{2.31}
\end{equation*}
$$

and finally

$$
\begin{equation*}
\Phi \sim(\mathrm{d} z / \mathrm{d} \rho)^{-1 / 2} \Phi_{0} . \tag{2.32}
\end{equation*}
$$

This approximation was justified by the methods of Lynn and Keller (1970) in the preceding paper.

The approximation (2.32) with $\Phi_{0}$ given by (2.31) is identical to our previous one, but with the difference that we have a totally different representation for $z$ as a function of $\rho$ (equation (2.19) as opposed to equation (2.3)).

The solution of (2.31) is (Arnold 1980)

$$
\begin{equation*}
\Phi_{0}=z^{m+1 / 2} \exp \left(-V z^{2} / 2\right) L_{\nu}^{(m)}\left(V z^{2}\right), \tag{2.33}
\end{equation*}
$$

and this solution satisfies the boundary condition (i) at the origin, as $z$ and $\rho$ are analytic functions of each other and $z \rightarrow 0$ as $\rho \rightarrow 0$. The parameter $\nu$ is defined by

$$
\begin{equation*}
2(2 \nu+m+1)=V z_{1}^{2} \tag{2.34}
\end{equation*}
$$

A contour integral for $L_{\nu}^{(m)}(u)$ is (Arnold 1977 $\dagger, 1980$ )

$$
\begin{equation*}
2^{m} \mathrm{e}^{-u / 2} L_{\nu}^{(m)}(u)=\frac{1}{2 \pi \mathrm{i}}\left(\mathrm{e}^{-\nu \pi \mathrm{i}} W_{1}+\mathrm{e}^{\nu \pi \mathrm{i}} W_{2}\right) \tag{2.35}
\end{equation*}
$$

with, for $j=1,2$,

$$
\begin{equation*}
W_{j}=\int_{C_{i}}\left(\frac{1+s}{1-s}\right)^{\nu+(m+1) / 2}\left(1-s^{2}\right)^{(m-1) / 2} \mathrm{e}^{-s u / 2} \mathrm{~d} s \tag{2.36}
\end{equation*}
$$

The contour $C_{1}$ passes from $s=-1$ to $s=\infty \mathrm{e}^{\mathrm{i} \sigma}(\sigma>0)$ and $C_{2}$ is the image of $C_{1}$ in the real axis, in the reverse direction.

Thus, we have established the same approximation as in the preceding paper, but with a different representation for $z$, facilitated by the assumed smallness of $\lambda$.

## 3. Non-uniform asymptotic solution

We have yet to determine the eigenvalue $W^{2}$, which is a parameter in $\nu$ through equations (2.11) and (2.34). This is determined through the boundary condition (1.4b) at $\rho=1$, and to do this we seek a simpler approximation than (2.32); we approximate $\Phi_{0}$ further in the vicinity of $\rho=1$.

Let $z \rightarrow z_{0}$ as $\rho \rightarrow 1$. Then, from equation (2.21),

$$
\begin{equation*}
\left[z_{0}^{2}-z_{1}^{2}+\mathrm{O}\left(\lambda^{3}\right)\right](\mathrm{d} z / \mathrm{d} \rho)^{2}=-\lambda \tag{3.1}
\end{equation*}
$$

(since $Q^{2}=0$ at $\rho=1$ ). This indicates that

$$
\begin{equation*}
z_{0}^{2}-z_{1}^{2} \sim \mathrm{O}(\lambda) \tag{3.2}
\end{equation*}
$$

Thus, as $\lambda \rightarrow 0\left(W^{2} \rightarrow 0\right), z_{0} \rightarrow z_{1}$. This is the confluence of boundary and caustic referred to in § 1 , and it becomes clear that the point $z_{0}$ will lie in the vicinity of $z_{1}$ if $\lambda$ is small enough. The Liouville-Green approximation to $\Phi_{0}$ will fail in this case; in fact we might expect the correct approximation to be in terms of Airy functions. The Liouville transform of equation (2.31),

$$
\begin{align*}
& \Phi_{0}=(\mathrm{d} \tau / \mathrm{d} z)^{-1 / 2} \Psi  \tag{3.3a}\\
& -\tau(\mathrm{d} \tau / \mathrm{d} z)^{2}=z_{1}^{2}-z^{2} \tag{3.3b}
\end{align*}
$$

[^0]would lead us to expect an asymptotic approximation for $\Psi$ (see Olver 1974),
\[

$$
\begin{equation*}
\Psi \sim \Psi_{0} \tag{3.4a}
\end{equation*}
$$

\]

where

$$
\begin{equation*}
\Psi_{0}=A_{1} \operatorname{Ai}\left(V^{2 / 3} \tau\right)+A_{2} \operatorname{Bi}\left(V^{2 / 3} \tau\right) \tag{3.4b}
\end{equation*}
$$

and $\mathrm{Ai}(X)$ and $\mathrm{Bi}(X)$ are the standard Airy functions (Abramowitz and Stegun 1965, Olver 1974). $A_{1}$ and $A_{2}$ are constants which we must find. To do this we note that (3.3) and (3.4) imply

$$
\begin{equation*}
\Phi_{0} \sim(\mathrm{~d} \tau / \mathrm{d} z)^{-1 / 2}\left(A_{1} \mathrm{Ai}\left(V^{2 / 3} \tau\right)+A_{2} \operatorname{Bi}\left(V^{2 / 3} \tau\right)\right) \tag{3.5}
\end{equation*}
$$

and we try to extract such an approximation from the integral representation (2.36). As this is precisely the method used by Slater (1960) to approximate the Whitaker function, which is essentially the same as our function $\Phi_{0}$, we need only state the results. By integrating (2.36) by the method of steepest descent (Chester et al 1957), expressing the result as a sum of Airy functions and comparing with (3.5), we obtain

$$
\begin{equation*}
A_{1}=A_{0} \cos (\nu \pi) \tag{3.6a}
\end{equation*}
$$

and

$$
\begin{equation*}
A_{2}=-A_{0} \sin (\nu \pi) \tag{3.6b}
\end{equation*}
$$

The constant $A_{0}$ may be set to unity because of the homogeneity of the differential equations.

By using (3.3a) and (3.4a) in the boundary condition (1.4b), we can obtain the eigenvalue equation

$$
\begin{equation*}
\mathrm{d} \Psi_{0} / \mathrm{d} \tau=K_{0} \Psi_{0}, \quad \tau=\tau_{0}, \tag{3.7}
\end{equation*}
$$

where

$$
\begin{equation*}
K_{0}=\left(\frac{\mathrm{d} \tau}{\mathrm{~d} z} \frac{\mathrm{~d} z}{\mathrm{~d} \rho}\right)^{-1}\left[K+\frac{1}{2}\left(\frac{\mathrm{~d} z}{\mathrm{~d} \rho}\right)^{-1} \frac{\mathrm{~d}^{2} z}{\mathrm{~d} \rho^{2}}+\frac{1}{2}\left(\frac{\mathrm{~d} z}{\mathrm{~d} \rho}\right)\left(\frac{\mathrm{d} \tau}{\mathrm{~d} z}\right)^{-1} \frac{\mathrm{~d}^{2} \tau}{\mathrm{~d} z^{2}}\right] \tag{3.8}
\end{equation*}
$$

and

$$
\begin{equation*}
\tau_{0}=\lim _{\rho \rightarrow 1}(\tau) \tag{3.9}
\end{equation*}
$$

Expressing $\Psi_{0}$ in terms of the Airy functions by equation (3.4b) and rearranging leads to

$$
\begin{equation*}
\tan (\nu \pi)=\frac{\operatorname{Ai}^{\prime}\left(V^{2 / 3} \tau_{0}\right)-V^{-2 / 3} K_{0} \operatorname{Ai}\left(V^{2 / 3} \tau_{0}\right)}{\operatorname{Bi}^{\prime}\left(V^{2 / 3} \tau_{0}\right)-V^{-2 / 3} K_{0} \operatorname{Bi}\left(V^{2 / 3} \tau_{0}\right)} \tag{3.10}
\end{equation*}
$$

where the prime implies differentiation with respect to the argument of the Airy functions and we have used (3.6) to evaluate the left-hand side. Equation (3.10) is to be solved in conjunction with equation (2.34) for $\nu$ (or for $z_{1}$, since equation (2.34) connects $\nu$ and $z_{1}$ directly).

In order to solve equation (3.10) asymptotically as $V \rightarrow \infty$, we need to estimate some of the parameters appearing in it. We already have equation (2.22) for $z_{1}$; in addition we require expressions for $z_{0}, \tau_{0}$ and various derivatives of $\tau$ and $z$ evaluated at $\rho=1$. The principal method for finding these values is to allow $\rho \rightarrow 1, \zeta \rightarrow \zeta_{0}$ in (2.8), (2.12),
(2.5) and (3.3b) or their differentiated forms. In this way we obtain the following expressions:

$$
\begin{align*}
& z_{0}^{2}=\alpha_{0}+\lambda\left(\alpha_{1}+\beta_{1}\right)+O\left(\lambda^{2}\right)  \tag{3.11}\\
& z_{0}^{2}-z_{1}^{2}=\lambda \beta_{1}+\mathrm{O}\left(\lambda^{2}\right)  \tag{3.12}\\
& \tau_{0}=\lambda\left(2^{-2 / 3} \alpha_{0}^{-1 / 3} \beta_{1}\right)+\mathrm{O}\left(\lambda^{2}\right) \tag{3.13}
\end{align*}
$$

where $\alpha_{0}$ and $\alpha_{1}$ are given by (2.12a) and (2.15) respectively, and

$$
\begin{equation*}
\beta_{1}=2^{2 / 3} \alpha_{0}^{1 / 3} \lim _{\rho \rightarrow 1}\left[\left(-\frac{\mathrm{d} Q^{2}}{\mathrm{~d} \rho}\right)^{-2 / 3}\right] \tag{3.14}
\end{equation*}
$$

The derivatives that we require are

$$
\begin{align*}
\left(\frac{\mathrm{d} \tau}{\mathrm{~d} z} \frac{\mathrm{~d} z}{\mathrm{~d} \rho}\right)^{-1} & =\left(\frac{\mathrm{d} \tau}{\mathrm{~d} \rho}\right)^{-1}=\left(\frac{\tau_{0}}{\lambda}+\mathrm{O}\left(\lambda^{2}\right)\right)^{1 / 2}  \tag{3.15}\\
& =\left(2^{-2 / 3} \alpha_{0}^{-1 / 3} \beta_{1}\right)^{1 / 2}+\mathrm{O}(\lambda) \tag{3.16}
\end{align*}
$$

Further expressions for the other derivatives in equation (3.8) can be obtained, but we shall show that they are in any case negligible to the order to which we are working.

We are now in a position to explain the purpose of the condition $\lambda \sim \mathrm{O}\left(V^{-2 / 3}\right)$. When this holds, the argument of the Airy functions in equation (3.10) is

$$
\begin{equation*}
V^{2 / 3} \tau_{0} \sim \mathrm{O}(1) \tag{3.17}
\end{equation*}
$$

by equation (3.13). If $\lambda$ is allowed to be larger than $\mathrm{O}\left(V^{-2 / 3}\right)$, the Airy functions may be replaced by their asymptotic expansions for large argument, in which case the righthand side of (3.10) is exponentially small, and we are returned to the case we considered earlier (Arnold 1980). The above condition on $\lambda$ is therefore a means of ensuring that the eigenvalue $W^{2}$ is small enough to justify separate treatment.

We now estimate the order of magnitude of the term $K_{0}$ in (3.10), given by equation (3.8). We have

$$
K=\frac{1}{2}+W \mathbf{K}_{m}^{\prime}(W) / \mathbf{K}_{m}(W)
$$

and this form invites us to consider two possibilities:
(i) $W$ is large, even though $\lambda$ is small;
(ii) $W$ is small.

Possibility (i) above can arise since

$$
\begin{equation*}
\lambda=W^{2} / V^{2} \sim \mathrm{O}\left(V^{-2 / 3}\right) \tag{3.18a}
\end{equation*}
$$

and therefore

$$
\begin{equation*}
W \sim \mathrm{O}\left(V^{2 / 3}\right) \tag{3.18b}
\end{equation*}
$$

Possibility (ii) may arise because it is obviously conceivable that, for some $V, W^{2}=0$. This case corresponds to cut-off of the propagating mode, and does not violate (3.18a).

If (i) above is true, then

$$
\begin{align*}
K & \sim-W+\frac{1}{2}+\mathrm{O}\left(W^{-1}\right)  \tag{3.19}\\
& \sim \mathrm{O}\left(V^{2 / 3}\right) \tag{3.20}
\end{align*}
$$

Therefore, as $V \rightarrow \infty$, the terms containing second derivatives in equation (3.8) are negligible, and we obtain

$$
\begin{equation*}
V^{-2 / 3} K_{0} \sim \mathrm{O}(1) \tag{3.21}
\end{equation*}
$$

In fact, if we let

$$
\begin{equation*}
X=V^{2 / 3} \tau_{0} \tag{3.22}
\end{equation*}
$$

then equation (3.10) may be approximated by

$$
\begin{equation*}
\tan (\nu \pi) \sim \frac{\operatorname{Ai}^{\prime}(X)+X^{1 / 2} \operatorname{Ai}(X)}{\operatorname{Bi}^{\prime}(X)+X^{1 / 2} \operatorname{Bi}(X)} \tag{3.23}
\end{equation*}
$$

when the precise form of $K_{0}$ is calculated. The value of $\lambda$ satisfying the boundary condition at $\rho=1$ is obtained by solving equations (3.23) and (2.34) for $\lambda$, by eliminating $\nu$.

We postulate the existence of an expansion

$$
\begin{equation*}
\lambda=V^{-2 / 3}\left[\lambda_{0}+\lambda_{1} V^{-1 / 3}+\mathrm{O}\left(V^{-2 / 3}\right)\right] . \tag{3.24}
\end{equation*}
$$

It turns out that $\lambda_{0}$ and $\lambda_{1}$ cannot be entirely independent of $V$, but are functions of the parameter

$$
\begin{equation*}
p=V^{-1 / 3} \Delta \tag{3.25}
\end{equation*}
$$

where

$$
\begin{align*}
& \Delta=V-2 N / \alpha_{0},  \tag{3.26a}\\
& N=2 q+m+1, \quad q \in\{0,1,2 \ldots\} . \tag{3.26b}
\end{align*}
$$

Nevertheless, since $V$ and $N$ can be regarded as independent of each other, we may also consider $p$ and $V$ to be mutually independent for the purpose of obtaining an asymptotic approximation for $\lambda$.

Since equation (3.23) is not changed if an integer $q$ is subtracted from $\nu$, we let

$$
\begin{equation*}
\nu=q+\theta, \quad 0 \leqslant \theta<1, \tag{3.27}
\end{equation*}
$$

and, in the same spirit as equation (3.24), we suppose that

$$
\begin{equation*}
\theta=\theta_{0}+\theta_{1} V^{-1 / 3}+\mathrm{O}\left(V^{-2 / 3}\right) \tag{3.28}
\end{equation*}
$$

and

$$
\begin{equation*}
X=X_{0}+X_{1} V^{-1 / 3}+\mathrm{O}\left(V^{-2 / 3}\right) \tag{3.29}
\end{equation*}
$$

Then substitution of (3.27), (3.28), (3.29), (3.24) in (3.23) and (2.34), (using (2.22) to express $z_{1}^{2}$ as a series in powers of $V^{-1 / 3}$ like (3.24)) and elimination of $\nu$ by equating similar terms leads to

$$
\begin{equation*}
\lambda_{0}=-\left(\frac{\alpha_{0}}{\alpha_{1}}\right) p=-\left(\frac{\alpha_{0}}{\alpha_{1}}\right) V^{-1 / 3} \Delta \tag{3.30}
\end{equation*}
$$

and

$$
\begin{equation*}
\lambda_{1}=\frac{1}{4 \pi \alpha_{1}} \tan ^{-1}\left(\frac{\mathrm{Ai}^{\prime}\left(X_{0}\right)+X_{0}^{1 / 2} \mathrm{Ai}\left(X_{0}\right)}{\operatorname{Bi}^{\prime}\left(X_{0}\right)+X_{0}^{1 / 2} \operatorname{Bi}\left(X_{0}\right)}\right) \tag{3.31}
\end{equation*}
$$

where

$$
\begin{equation*}
X_{0}=2^{-2 / 3} \alpha_{0}^{-1 / 3} \beta_{1} \lambda_{0} . \tag{3.32}
\end{equation*}
$$

Equations (3.30) and (3.31) are the final expressions for the coefficients in the expansion of $\lambda$, and, since $\lambda=W^{2} / V^{2}$, we may compute $W^{2}$ from these equations. The parameters $\alpha_{0}, \alpha_{1}$ and $\beta_{1}$ are given by (2.12a), (2.15) and (3.14) respectively. (It should be noted that, despite the minus sign in (3.30), $\lambda_{0}$ is positive because $\alpha_{1}$ is negative.) The most useful expression we can derive is for $W^{2} / V$ :

$$
\begin{equation*}
W^{2} / V=-\alpha_{0} \Delta / \alpha_{1}+\lambda_{1}+\mathrm{O}\left(V^{-1 / 3}\right) \tag{3.33}
\end{equation*}
$$

with $\lambda_{1}$ given by (3.31) and (3.32).
To obtain some insight into this expression, it is instructive to consider the case of exactly quadratic index variation,

$$
\begin{equation*}
f=\rho^{2}, \quad Q^{2}=1-\rho^{2} . \tag{3.34}
\end{equation*}
$$

Then

$$
\begin{equation*}
\alpha_{0}=-\alpha_{1}=1, \tag{3.35}
\end{equation*}
$$

and we have

$$
\begin{equation*}
W^{2} / V=\Delta+\lambda_{1} . \tag{3.36}
\end{equation*}
$$

The equivalent result neglecting the effect of the cladding medium, which we have shown elsewhere is accurate for modes not close to cut-off (Arnold 1980), would be

$$
\begin{align*}
W^{2} / V=V-U^{2} / V & =V-2 N  \tag{3.37}\\
& =\Delta \tag{3.38}
\end{align*}
$$

This clearly shows that the term $\lambda_{1}$ in (3.33) is due to the presence of the finite cladding boundary. As $\Delta$ becomes sufficiently large, with $V$ fixed, $\lambda_{0}$ and $X_{0}$ also have large values, and the term $\lambda_{1}$ becomes exponentially small through the behaviour of the Airy functions in equation (3.31). Since large $\Delta$ corresponds to large $W^{2}$, this is in exact agreement with our previous result (Arnold 1980), that for modes not close to cut-off the effect of the finite boundary is an exponentially small correction to the eigenvalue for an unbounded medium.

This completes our analysis of the case (i) above. We now repeat the analysis assuming (ii) above is true; $W$ may not be large enough for (3.19) to hold.

In that case we find that

$$
\begin{equation*}
K \sim \mathrm{O}(1) \tag{3.39}
\end{equation*}
$$

and so

$$
\begin{equation*}
V^{-2 / 3} K_{0} \sim \mathrm{O}\left(V^{-2 / 3}\right) \tag{3.40}
\end{equation*}
$$

Then it is clear that (3.10) may be approximated by

$$
\begin{equation*}
\tan (\nu \pi) \sim \frac{\operatorname{Ai}^{\prime}\left(V^{2 / 3} \tau_{0}\right)}{\operatorname{Bi}^{\prime}\left(V^{2 / 3} \tau_{0}\right)}+\mathrm{O}\left(V^{-2 / 3}\right) \tag{3.41}
\end{equation*}
$$

Since none of the other approximations made in case (i) are affected, we may carry the analysis through to obtain a modified expression for $\lambda_{1}$ :

$$
\begin{equation*}
\lambda_{1}=\frac{1}{4 \pi \alpha_{1}} \tan ^{-1}\left(\frac{\operatorname{Ai}^{\prime}\left(X_{0}\right)}{\operatorname{Bi}^{\prime}\left(X_{0}\right)}\right) . \tag{3.42}
\end{equation*}
$$

Equations (3.30) and (3.32) remain valid as written.

It is of some interest to calculate the value of $V$ for which $W=0$ exactly. In that case we find

$$
\begin{equation*}
K_{0} \sim \mathrm{O}(1) \tag{3.43}
\end{equation*}
$$

and

$$
\begin{equation*}
\tau_{0}=0 \tag{3.44}
\end{equation*}
$$

and so equation (3.10) becomes

$$
\begin{align*}
\tan (\nu \pi) & =\frac{\mathrm{Ai}^{\prime}(0)}{\operatorname{Bi}^{\prime}(0)}+\mathrm{O}\left(V^{-2 / 3}\right)  \tag{3.45}\\
& =-3^{-1 / 2}+\mathrm{O}\left(V^{-2 / 3}\right) . \tag{3.46}
\end{align*}
$$

Hence

$$
\begin{equation*}
\nu=q-\frac{1}{6}+\mathrm{O}\left(V^{-2 / 3}\right) \tag{3.47}
\end{equation*}
$$

Also we have, since $\lambda=0$,

$$
\begin{equation*}
z_{1}^{2}=\alpha_{0} \tag{3.48}
\end{equation*}
$$

and so, equation (2.34) is equivalent to

$$
\begin{equation*}
2 N-\frac{2}{3}+\mathrm{O}\left(V^{-2 / 3}\right)=V \alpha_{0} \tag{3.49}
\end{equation*}
$$

and finally

$$
\begin{equation*}
V=\left(2 N-\frac{2}{3}\right) / \alpha_{0}+\mathrm{O}\left(V^{-2 / 3}\right) \tag{3.50}
\end{equation*}
$$

is the final expression for the cut-off frequency, at which $W=0$. When the index variation is exactly quadratic,

$$
\begin{equation*}
\alpha_{0}=1 \tag{3.51}
\end{equation*}
$$

and (3.50) agrees with the correct result for this case (Arnold 1977). Since the error term in (3.50) is $\mathrm{O}\left(V^{-2 / 3}\right)$, this approximation is only valid for a mode whose cut-off frequency is high enough to ensure the smallness of this term.

## 4. Conclusions

It has here been shown how approximate expressions can be obtained for an eigenvalue problem arising in the theory of inhomogeneous dielectric waveguides; in particular, we have been concerned with the case of those modes which are close to cut-off, where the existence of finite boundaries has a significant effect on the asymptotic ( $V \rightarrow \infty$ ) behaviour of the eigenvalues. The use of uniform asymptotic approximations to the solution has been demonstrated, and the further approximations necessary to solve the eigenvalue problem have been evaluated. It is hoped that these methods may eventually lead to a complete and systematic theory of propagation in waveguides having arbitrary variation of refractive index, under the assumption $V \rightarrow \infty$. Detailed calculations using these expressions are to be described elsewhere.

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[^0]:    $\dagger$ Several typographical errors occurred in this paper; see reference.

