Non-Canonical Scalar Fields and Their Applications in Cosmology and Astrophysics

by

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To my Father, my first science teacher
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ABSTRACT

Non-Canonical Scalar Fields and Their Applications in Cosmology and Astrophysics

by

Christopher S. Gauthier

Chair: Professor Ratindranath Akhoury

In this thesis we will discuss several issues concerning cosmological applications of non-canonical scalar fields, which are generically referred to as $k$-essence. First, we consider two examples of $k$-essence. These are the rolling tachyon and static spherically symmetric solutions of non-canonical scalar fields in flat space. We find constraints on the form of the allowed interactions in the first case and on the choice of boundary conditions in the latter. For the rolling tachyon we find that at late times the tachyon matter behaves like a non-relativistic dust, thus making it a dark matter candidate. For the static spherically symmetric solutions we show that solutions which are finite at the origin must have negative energy density there.

Next, we consider static spherically symmetric solutions of non-canonical scalar fields coupled to gravity as a way to explain dark matter halos as a coherent state of the scalar field. Consistent solutions are found with a smooth scalar profile which can describe observed rotation curves. The non-trivial solutions have negative energy density near the origin, though the total energy is positive. We also reconsider the no scalar hair theorems for black holes with emphasis on asymptotic boundary conditions and superluminal propagation.
After this we show that, for general scalar fields, stationary configurations are possible for shift symmetric theories only. This symmetry with respect to constant translations in field space should either be manifest in the original field variables or reveal itself after an appropriate field redefinition. In particular this result implies that neither k-essence nor quintessence can have exact steady state/Bondi accretion onto black holes. Finally, we find that stationary field configurations are necessarily linear in Killing time, provided that shift symmetry is realized in terms of these field variables.

The next discussion outlines a general program for reconstructing the action of non-canonical single field inflation models from CMBR power spectrum data. This method assumes that an action depends on a set of undetermined functions, each of which is a function of either the inflaton field or its time derivative. The scalar, tensor and non-gaussianity of the curvature perturbation spectrum are used to derive a set of reconstruction equations whose solution set can specify up to three of the undetermined functions. This method is used to find the undetermined functions in various types of actions assuming power law type scalar and tensor spectra.

Finally, we study a novel means of coupling neutrinos to a Lorentz violating k-essence background. We first look into the effect k-essence has on the neutrino dispersion relation, and derive the neutrino velocity in a k-essence background. Next, we look at the effect on neutrino oscillations. It is found that if k-essence couples non-diagonally to the neutrino flavor eigenstates, this leads to an oscillation length that varies with the neutrino energy like $\lambda \sim E^{-1}$. This is to be compared with the $\lambda \sim E$ dependence seen in mass-induced neutrino oscillations. While such a scenario is not favored experimentally, it places tight constraints on the possible interaction that a k-essence background can have with neutrinos.
CHAPTER I

Introduction

The past two decades have been an exciting time for cosmologists and astrophysicists. This period has seen the beginning of true, precision cosmology. What was once a field that relied almost exclusively on qualitative predictions, was now able to test theories with definitive numerical accuracy. In this era we have seen two satellite missions: the Cosmic Background Explorer (COBE) and the Wilkinson Microwave Anisotropy Probe (WMAP), measure the spectrum of fluctuations in the cosmic microwave background radiation (CMBR). The first of these missions, COBE, confirmed that the CMBR has a black body spectrum [6], which was the final nail in the coffin for all competitors to the Big Bang Theory at that time. Both COBE and WMAP have also measured the minute temperature variations in the CMBR, leading to widespread support for cosmological inflation [7, 8].

During the period between COBE and WMAP, astrophysicists were able to measure the rate of the universe’s expansion. Two teams of astrophysicists [9] used type IA supernovae as “standard candles”, whose absolute brightness are remarkably consistent, and thus their distance away from us can be determined independently of their red-shift. This in turn allowed them to measure the rate of the expansion, which to the surprise of many, turned out to be increasing with time.
As exciting and important as these discoveries were for cosmology, they raised more questions than answers. Although COBE and WMAP confirmed many of the predictions of cosmological inflation, they could not tell us conclusively what mechanism was responsible for inflation. Similarly, while we have determined that the expansion rate of the universe has (at least recently) been accelerating, we still have no idea why it is accelerating.

1. Inflation

Up until the early 1990’s, the status of the Big Bang model was still in doubt. Although it was the favored model by that time, there were still problems with the theory that gave some physicists reason to object. In particular, there were three outstanding problems with the Big Bang. The first of these was called the horizon problem. In the classical model of the Big Bang, the comoving causal horizon at the time of photon decoupling was approximately $180\Omega_0^{-1/2}h^{-1}$ Mpc. However, the present comoving horizon of the universe is considerably larger; $5820h^{-1}$ Mpc$^1$. This means that our observable universe today consists of approximately $10^5$ regions that were causally disconnected from each other at the time when the CMBR settled into its present state. However, despite this, the CMBR is remarkably smooth across the entire sky. Such homogeneity over causally disconnected regions is something the classical Big Bang can not explain satisfactorily.

The second problem for the classical Big Bang model is the flatness of the universe’s geometry. If one does a simple study of the Friedmann equations, they would quickly raise the issue that a flat universe is not a stable solution. If we write the 

\footnote{This assumes that the universe is flat}
Friedmann equations in terms of the density parameter $\Omega$ what we find is that

\[
\frac{K}{a^2 H^2} = \Omega - 1,
\]

(1.1)

where $\Omega = \frac{\rho}{3M^2_{pl} H^2}$, and $H = \frac{\dot{a}}{a}$ is the Hubble parameter. Here, $K$ is a constant that represents the intrinsic spatial curvature of the universe, and the scale factor $a$ can always be defined such that $K = 0, \pm 1$. If $K$ vanishes then it follows that the density parameter $\Omega$ is a constant and equal to the critical density $\Omega = 1$. When this happens the universe has a flat, Euclidean spatial geometry. On the other hand, if $K$ is nonzero then $\Omega$ is not at the critical value. If the expansion of the universe is decelerating, then the factor $a^2 H^2 = \dot{a}^2$ is getting smaller, which means that the deviation of $\Omega$ from the critical density will continue to get larger as time goes on.

This was a troubling notion for supporters of the Big Bang at that time because all evidence seemed to indicate that we live in an extremely flat universe. Current estimates place the present day value of $|\Omega - 1|$ at less than 0.01. This means that our universe is so flat that in order to end up within the presently observed range of $\Omega$, then at the time of the Planck era the deviation of $\Omega$ from the critical density would have to be smaller than one part in $10^{60}$.

Finally, the third big problem with the Big Bang is the paradoxically low density (if not the absence) of monopoles. If our theory of gauge fields is correct, which all Earth based experiments suggest, then a copious number of monopoles should have been produced at the beginning of the Big Bang [10]. This absence (so far) of monopoles in the universe is another puzzle that the classical Big Bang is unable to account for.

As luck would have it, Alan Guth suggested a simple solution to these problems [11], a solution that would later be known as cosmological inflation (or simply in-
flation). Guth’s model involved giving the universe a short period of accelerated expansion. This rapid expansion is powered by a positive vacuum energy density that is sourced by a scalar field, called the inflaton, that lies in a metastable vacuum. Since this state is unstable due to quantum fluctuations, the inflaton will eventually tunnel out of the metastable vacuum and into a stable one, creating bubbles of true vacuum in the bulk of false vacuum. If these bubbles are produced at a large enough rate, they will collide with one another, releasing the energy stored in the walls of the bubbles as radiation, thereby reheating the universe and ending inflation.

Guth’s initial motivation for this idea was to solve the monopole problem. In the inflationary universe, monopoles created just before the inflationary expansion are diluted away by the rapidly increasing volume of our observable universe. Guth’s model, however, turned out to solve more than just this problem. As it turns out an initial rapid expansion also solved the horizon and flatness problems. In regards to the horizon problem, inflation reconciled this paradox by allowing for the observable universe to occupy a casual patch before inflation started, during which the universe could achieve thermal equilibrium. Inflation also gave an explanation for the flatness problem since during an accelerated expansion $\dot{a}$ is getting larger, and thus as (1.1) implies, $\Omega$ will be driven towards the critical density.

Unfortunately, Guth’s original inflation model had a serious problem that made it unworkable; a fact he acknowledged when he first introduced his model. In Guth’s model of inflation, the decay rate of the metastable vacuum has to be small enough to provided for a sufficient amount of inflation, but has to be large enough for the bubbles of true vacuum to be produced at a fast enough rate for them to collide and reheat the universe. These requirements were found to be incompatible, leaving Guth’s model dead on arrival. However, the idea that an early exponential expansion
could solve three of the biggest objections to the Big Bang was too tantalizing to give up. Fortunately, soon after Guth published his findings, Andrei Linde [12] and independently Andreas Albrecht and Paul Steinhardt [13], came up with a new model of inflation that was soon called slow roll inflation. In this realization of inflation the inflaton is a scalar field with a potential that has a zero energy vacuum state. When the inflaton is displaced from the vacuum, the inflaton potential will be nonzero, providing the vacuum energy needed to drive inflation. When the inflaton reaches the vacuum, inflation ends and the energy in the inflaton field is converted into radiation that reheats the universe. Today most models of inflation are based on the paradigm set by Albrecht and Steinhardt, and Linde.

If all inflation did was explain away problems with the Big Bang, then inflation would be regarded by critics of the Big Bang as nothing more than an ad hoc fix. However, as it turns out inflation makes a testable prediction. Due to quantum fluctuations in the behavior of the inflaton, the value of the vacuum energy density will vary at different points in space. Through a careful analysis [14] it was shown that this variation, although slight, would eventually leave its mark as spatial fluctuations in the CMBR. Specifically, nearly all inflation models make the generic prediction that at large distance scales the magnitude of the fluctuations in the CMBR temperature, should be relatively independent of the size of the region of the fluctuation. The possibility of confirming this prediction was one of the reasons that the COBE satellite was created. COBE’s observation of nearly scale invariant temperature fluctuations solidified the support for cosmic inflation and put to rest any serious doubts about the Big Bang.
2. Dark Energy

After Edwin Hubble made his landmark observation that the universe was expanding, many physicists assumed that due to the universal gravitational attraction of matter, this expansion rate should be getting smaller as time goes on. However, in 1998 this view was quickly thrown out of the window [9]. As it turns out not only is the expansion of the universe not decreasing, but it is actually accelerating. At present, the best answer we can give for this accelerating expansion rate is the existence of a positive vacuum energy density pervading the universe. This vacuum energy can be modeled in general relativity as a cosmological constant term. The cosmological constant (c.c.) was originally proposed by Einstein as a means to create a model of a static universe. After Einstein heard about Hubble’s discovery, he recanted his belief in a static universe, and declared his inclusion of the c.c. “his biggest blunder”. However, with these recent insights it seems that Einstein’s blunder may not have been a blunder after all. Although the accelerated expansion can be modeled by the inclusion of a c.c., this does not explain why a c.c. exists.

For particle physicists, the question was not so much why there is a c.c., but why is it so small? In the calculation of loop corrections to the effective action of the standard model, the correction to the c.c. term $\Lambda$ is of the order of the Planck scale $\Lambda \sim M_{pl}^4 \sim 10^{109}$ eV$^4$. This is staggeringly higher than the currently estimated value of the c.c, which is $\Lambda \sim 10^{-11}$ eV$^4$; off from the expected value by a factor of $10^{120}$! This gigantic discrepancy has been dubbed the “the worst prediction in physics”. While the value of the corrections to the c.c. can be eliminated by counter terms during renormalization, it is troubling that such large numbers must be contrived in such a way as to result in a number that is smaller by a factor of $10^{-120}$ of the
original numbers.

Unfortunately, there has been little success in trying to explain the nature of dark energy. Some of the ideas that have been proposed include: *quintessence* \( [15, 16] \), *anthropic selection* \( [17] \), *extremely large cosmic voids* \( [18] \). One proposed explanation, which uses a scalar field with a non-canonical action as the source of the c.c., is called *k-essence* \( [19, 20, 21] \). This model has several novel features and will be a primary focus of this thesis.

3. *K-essence and K-inflation*

Dynamical scalar fields have held a special place in modern physics. When physicists encounter new phenomena, often times a phenomenological model based on scalar fields is used in order to gain some initial understanding. Examples of these include: the Yukawa model of the strong force \( [22] \), the Landau-Ginzburg model of superconductivity \( [23] \), the Higgs model \( [24] \), the axion \( [25] \), and most recently cosmological inflation \( [11, 12, 13] \). Scalar fields enjoy the advantage of having the simplest behavior under general space-time coordinate reparametrizations, making them natural tools for building toys models. The action of a real relativistic scalar field \( \phi \), in a general \( N \)-dimensional space-time background is given by

\[
S_{\text{scalar}} = \int d^N x \sqrt{-g} \left[ \frac{1}{2} g^{\mu \nu} \partial_\mu \phi \partial_\nu \phi - V(\phi) \right],
\]

where \( g_{\mu \nu} \) is the background space-time metric, \( g = \text{det}(g_{\mu \nu}) \), and \( V(\phi) \) is the potential energy function of \( \phi \). Of all field theory actions, this is by far the simplest. In addition to their simplicity, scalar fields also have the unique and useful property of being able to have a constant nonzero vacuum expectation value (i.e., \( \langle \phi \rangle \neq 0 \)) while still maintaining Lorentz invariance. This property is a prized feature of scalar fields, and it is the reason that models like the Higgs mechanism and inflation are
possible.

Although (1.2) is the most widely used theory of scalar fields, it would break no known fundamental physical principles if more general scalar field actions were considered, such as

$$S = \int d^N x \sqrt{-g} L(X, \phi),$$

where $X = \frac{1}{2} g^{\mu\nu} \partial_\mu \phi \partial_\nu \phi$. Here the function $L$ is the Lagrangian of the scalar field and can be almost any function of two variables. While there are no theoretical objections to such general scalar field actions, many physicists paid little attention to such possibilities, as there didn’t seem to be any good reason to consider them. Since there are an infinite number of possible functions $L$ to consider, without some motivation from theory or experiment, little can be done in the way of studying them.

Recently, however, this class of theories has started to attract attention due to new discoveries in both theory and experiment. On the theory side, advances in our understanding of string theory have led to a deeper appreciation of the importance of D-branes [26]. The effective action of the D-brane degrees of freedom can be cast into the form of the Dirac-Born-Infeld (DBI) action, an action that was previously proposed as an alternative to the standard electrodynamical action. The simplest DBI action describes the motion of a D3-brane along a single direction:

$$S_{DBI} = - \int d^4 x \sqrt{-g} \left[ f^{-1}(\phi) \sqrt{1 - 2f(\phi)X - f^{-1}(\phi) + V(\phi)} \right].$$

The features of this action will be discussed in depth in chapter §V. For now we should note that this is an example of a non-canonical scalar field action, and has found possible applications in cosmology [27, 28].
In the realm of phenomenology, some have considered the possibility of using non-canonical scalars to explain the coincidence problem [15, 29]. Put simply, the coincidence problem refers to the fact that the dark energy, which is now the main driver of cosmic expansion, has only recently been the dominant component of the total energy density of the universe. If this were not true, and dark energy had become dominant earlier on, it is inconceivable that matter would be able to form clusters large enough to create galaxies. In order to get galaxies there must have been a period of matter domination, between the periods of radiation and dark energy domination. Attempts have been made to explain this coincidence as a consequence of the behavior of dark energy. These proposals have concentrated on models that exhibit the so called tracker solution. In these kinds of models, a scalar field is the source of the dark energy, and its action is setup so that the equation of state of the scalar field only starts to behave like vacuum energy after matter has come to dominate over radiation.

One model that has tracking solutions is *quintessence* [15, 16]. Quintessence uses a canonical scalar field with a suitably chosen potential to get the correct tracking solution needed to solve the coincidence problem. While it is possible to solve the coincidence problem with quintessence, it was quickly realized that to do so requires an incredibly high level of fine tuning; defeating the entire purpose behind quintessence in first place.

An alternative to quintessence, dubbed *k-essence* (the “k” standing for kinetic), works the same way, but uses non-canonical scalar fields as the source of the vacuum energy [19, 20, 21]. The advantage k-essence has over quintessence is that tracker solutions are the general solutions of k-essence. Thus, k-essence does not require the high level of fine tuning that quintessence does if one wishes to solve the coincidence
problem. However, the naturalness of k-essence does come at a price. It was discovered [30] that any k-essence model created to solve the coincidence problem would have a brief period when fluctuations in the k-essence field propagate at superluminal speeds. At first this may seem like a deal breaker for k-essence; however, it was later shown by the authors in [31] that despite the superluminal speeds, causality is preserved in k-essence.

An interesting feature of k-essence is that, unlike quintessence, the kinetic energy of k-essence can source the cosmological constant. Vacuum energy has the critical property that it is constant with respect to the scale factor $a$ of the universe. The kinetic energy of a canonical scalar field leads to an energy density that goes like $\propto a^{-6}$, which means that it dissipates too quickly, even more quickly than matter ($\propto a^{-3}$) and radiation ($\propto a^{-4}$). However, because of the non-trivial dependence of the k-essence Lagrangian on $X$, it is possible to get a vacuum energy entirely through the kinetic energy. Therefore, it is possible to have so called “kinetic” k-essence theories; actions that only depend on $X$.

This ability to create vacuum energy has also made k-essence a possible candidate for the inflaton. Theories that attempt to explain inflation in the context of k-essence have been dubbed $k$-inflation [32]. K-inflation models have a number of interesting features that set them apart from the typical slow roll inflation models. One such feature of k-inflation is that fluctuations in the inflaton field can propagate at speeds different than the speed of light. When viewed as a continuous classical medium, the speed of fluctuations $c_s$ in a general scalar field is

$$c_s^2 = \frac{\partial p/\partial X}{\partial \rho/\partial X} = \frac{L_X}{L_X + 2XL_{XX}}. \tag{1.3}$$

Here $p$ and $\rho$ denote the pressure and energy density of the scalar field. The speed $c_s$ is typically referred to as the sound speed. Clearly, from (1.3) we can see that
$c_s$ will be equal to one for a canonical scalar field$^3$. However, for k-essence $c_s$ will in general be different from, and can even exceed, the speed of light [31, 33]. This can lead to some interesting consequences in the case of k-inflation models since the spectrum of primordial curvature perturbations depends on the sound speed of the k-essence perturbations [34].

Finally, another reason k-inflation has been under consideration is because they often predict large non-gaussianities [35, 36]. In canonical inflation models, non-gaussianities can only be produced through interactions that are cubic or higher order in the inflaton field variable, or indirectly through the inflaton’s interaction with gravity. However, the non-gaussianity produced in these ways is on the same order as the slow roll parameters [37, 38], and is therefore small. Non-gaussianities in k-inflation on the other hand can be quite substantial due to the possible non-linear dependence of the action on $X$. Thus, non-gaussianities are an important discriminator between canonical and non-canonical inflation models. The application of non-gaussianities to finding the form of the k-inflation action is discussed in chapter §V of this thesis.

Therefore, non-canonical scalars have the ability to explain the late-time accelerated expansion of the universe, and the related coincidence problem through k-essence; and they can also provide an alternative explanation for the primordial accelerated expansion of the universe through k-inflation.

4. Dark Matter

The idea that the universe contains particles that we have yet to observe is an old idea in physics. Several times when physicists were faced with an unexpected observation or a difficulty with theoretical models, a new, as yet unseen particle

$^3$Throughout this thesis we will take the natural units so that $c = 1$. 

11
turned out to be the solution. This has happened in the case of the neutrino [39],
the $\pi$ meson [22], and the charm quark [40]. In light of this trend, it’s no surprise that
when galactic rotation curves were found to be in conflict with the naive expectation
of Newtonian mechanics [41], physicists immediately began to propose all kinds of
new particles to explain the discrepancy. According to the results from the surveys of
galactic rotation curves, the distribution of observable matter in all galaxies studied
so far indicate that there is not enough observable matter on the edges of these
galaxies to explain the high rotational velocity there. In light of these observations,
it is reasonable to suggest that there is additional matter that, although invisible to
our telescopes, nevertheless comprises the majority of the mass content of galaxies,
and by extension our universe.

The history of dark matter (so named because it is required by definition to be
decoupled from electromagnetic interactions) has a long history in astrophysics, and
yet in all this time, its true nature has so far eluded our understanding. At first it
was suggested that dark matter may have a much more mundane original; perhaps
composed of the remnants of stars, or the puny sized gas giants that weren’t able
to become stars. For a time, neutrinos were considered a possible candidate for
dark matter [42]. However, as we have learned more about the required properties
of dark matter, it has become apparent that these run of mill solutions are not
enough to explain the observations. Clearly more unconventional options have to be
considered. Many of the current explanations for dark matter rely on the existence
of some exotic, as yet unseen, particle. The proposed identity of the dark matter
candidate has been primarily motivated by theories that go beyond the standard
model, such as supersymmetry [43]. Computer simulations of structure formation in
the early universe have also helped in determining the required properties of dark
As popular as this approach is among physicists, there are other, albeit more radical alternatives to dark matter. One of these is *Modified Newtonian Dynamics* (MOND), which attempts to explain the observed discrepancies as a modification of gravity at weak gravitational fields [44]. These models, however, have met with limited success. Another proposal for explaining dark matter, that we will be studying in chapters §II-IV, is the idea that dark matter is really a k-essence condensate. In these chapters we will discuss the conditions needed to have stable k-essence halos around galaxies.

5. Neutrino Oscillations

If there is a k-essence field pervading the universe, then it is natural to ask what effects it might have on the standard model family of particles. From the perspective of effective field theory, it is almost a given that if a k-essence background exists, it will directly couple to all other matter that is present. Due to its importance in cosmology and astrophysics, the neutrino presents the most interesting case in which to study possible k-essence/matter interactions.

In 1930 Wolfgang Pauli [39] proposed the existence of a (at that time undetected) particle, later called the neutrino, that was responsible for carrying away the energy and momentum that was missing in the decay products of beta decay. Observations of the energy spectra of the beta decay products led to the initial conclusion that the neutrino was massless. However, the notion that neutrinos may have a small but nonzero mass has been a persistent idea in modern particle physics. For some time the data did not imply any need for a nonzero neutrino mass, until neutrino observatories measured a deficit in the expected number of neutrinos being produced.
by our sun [45]. Even before this anomaly was found physicists had considered the possibility of mass causing just this kind of result [46]. In order to understand how neutrino mass can explain the solar neutrino problem we first have to note that neutrinos come in three different varieties: the electron neutrino, the muon neutrino, and the tau neutrino. Early solar neutrino detectors were only sensitive to the electron neutrino species, since this is the only neutrino flavor that is produced in the kinds of reactions that take place in the sun.

It was found that if the various neutrino flavors have different masses, then this can lead to the phenomenon of neutrino oscillations. This means that if a neutrino of the electron variety is created at the source (in this case the sun) then as it travels to an Earth based detector, the probability of observing this neutrino in the electron neutrino state will oscillate with respect to the distance traveled. A simplified analysis of neutrino oscillations when there are $N$ different flavors, shows that the probability of observing an $f'$ flavored neutrino some distance $L$ away from the source of the neutrino where it was $f$ flavored is

$$P(\nu_f \to \nu_{f'}) = \delta_{f f'} - 4 \sum_{i>j} \text{Re}(U_{fi}^* U_{f'i} U_{f'j} U_{fj}^*) \sin^2 \left( \frac{\Delta m^2_{ij} L}{4E} \right)$$

(1.4)

$$+ 2 \sum_{i>j} \text{Im}(U_{fi}^* U_{f'i} U_{f'j} U_{fj}^*) \sin \left( \frac{\Delta m^2_{ij} L}{2E} \right).$$

Here, $E$ is the energy of the neutrino and $\Delta m^2_{ij} = m_i^2 - m_j^2$, where the $m_i$'s are the masses of the neutrino mass eigenstates. The indices $i$ and $j$ enumerate the neutrino mass eigenstates and assume the values $i, j \in \{1, 2, ..., N\}$. Here the $U_{fi}$'s are the matrix elements of the so called neutrino mass mixing matrix. Essentially, the matrix $U$ is an $SU(N)^4$ matrix that transforms the neutrino mass eigenstates into the neutrino flavor eigenstates. As long as $U$ is non-diagonal, then the last two

\footnote{This assumes that the neutrinos are Dirac fermions. If they are Majorana fermions then the mixing matrix must be an $U(N)$ matrix.}
terms of (1.4) will be nonzero and neutrinos will undergo flavor oscillations.

Thus, even if only electron neutrinos are being produced in the sun, as long as there is a mass difference between the neutrino mass eigenstates, some of those electron neutrinos will “flip” into mu or tau neutrinos when they arrive on Earth. Since the early neutrino observatories could not detect the mu and tau neutrinos, those electron neutrinos that flipped could not be detected. Therefore, the solar neutrino flux would have been lower than expected [47]. Today, these “lost” neutrinos have been detected, and are enough to account for the missing solar neutrino flux [48], confirming that neutrinos do indeed undergo flavor oscillations. This has been widely interpreted as evidence that neutrinos have mass\(^5\). However, other mechanisms that can also induce neutrino oscillations [49, 50, 51, 52, 53] are often ignored. Thus, it is not a given that neutrino oscillations confirm the existence of neutrino masses.

The critical feature needed to induce neutrino oscillations is a spatially varying phase difference between the energy eigenstates of the neutrino. Although mass can create this phase difference, there are many other ways in which to create flavor oscillations. Any coupling term in the neutrino Lagrangian that is non-diagonal in the neutrino flavor space is capable of inducing neutrino oscillations. In chapter §VI, we will study the possibility of k-essence induced neutrino oscillations (KINO) [5]. In this new mechanism, the neutrinos couple to k-essence by replacing the space-time metric in the neutrino action with a k-essence induced metric. In addition to modifying the dispersion relation of the neutrinos, we will show that if this coupling is non-diagonal in the neutrino flavor basis, neutrino oscillations occur even in the absence of neutrino masses.

\(^5\)To be more accurate, observation of the missing neutrino flux is evidence that the neutrino mass eigenstates have different masses. It is still possible for at least one of the neutrinos to be massless.
5. Outline of Thesis

This thesis is organized as follows. Chapters §II and §III discuss the possibility of dark matter halos as k-essence condensates. In chapter §II we investigate k-essence halos in flat space and in the FRW metric. In particular, we study the specific cases of the rolling tachyon, and static spherically symmetric solutions of general k-essence Lagrangians. In chapter §III we extend the analysis of static spherically symmetric solutions in chapter §II to include the effects of gravitational back-reaction. Chapter §IV is an outgrowth of the work done in the previous two chapters, and here the question of when stationary solutions are possible in k-essence is addressed. The existence of stationary k-essence solutions is addressed in the context of black hole accretion, where one typically uses such solutions in order to simplify the analysis. In chapter §V we change the subject from k-essence to k-inflation, and discuss a procedure for reconstructing the k-inflation action from cosmological observables. Next, in chapter §VI we explore the possibility of k-essence/neutrino interactions. In that chapter we will study the effects that a hypothetical k-essence/neutrino coupling has on the motion and flavor oscillations of neutrinos. Finally, in chapter §VII we review our major findings in the conclusion.
CHAPTER II

Classical Solutions of K-essence Theories

2.1 Introduction

Scalar field theories with higher derivatives play an essential role in the effective field theory approach (for reviews see [54]). An example is provided by the chiral Lagrangian which provides a good description of the strong dynamics at low energies. Applications of higher derivative theories to cosmology have also become popular in the last few years: examples here are effective field theories of the rolling tachyon [55, 56], DBI inflation [27, 28, 57], and k-essence [19], which attempts to provide a dynamical explanation of the so-called coincidence problem and the accelerated expansion of the universe. Recently [35], such higher derivative actions have been shown to enhance the non-gaussian fluctuations in the cosmic microwave background.

Theories with higher derivatives have the possibility of modifying the dispersion relations and hence may potentially lead to superluminal propagation. This aspect has been studied in detail in [58] where it was shown that causality and the absence of superluminal propagation require certain coefficients of the effective Lagrangian to be positive definite, which in turn has consequences for phenomenology [58, 59]. Thus, the constraints of causality and hyperbolicity of the equations of motion play a particularly important role in such theories. Another recent striking example is
the no-go theorem proved in [30]. It was shown there that in the context of the
original k-essence theories [19], it is impossible to simultaneously resolve the coincidence problem and the accelerated expansion of the universe without encountering problems with superluminal propagation.

In this chapter we apply the constraints [60, 61, 62, 63] that the equation of motion of the scalar field has a well defined initial value problem and there is no superluminal propagation of the small fluctuations around classical solutions in higher derivative theories. In particular we discuss in sections §2.3 and §2.4, respectively, the case of the rolling tachyon and the static solutions to the equations of motion of a general scalar theory with higher derivatives. For the case of the tachyon we consider a general Lagrangian of the form \( L = V(\phi)K(X) \), with \( V \) the potential of the tachyon and \( X = \frac{1}{2} g^{\mu\nu} \partial_\mu \phi \partial_\nu \phi \), and find the constraints on \( K(X) \) and the potential such that the energy density is finite but the equation of state parameter goes to zero at late times up to small corrections. We find that in order achieve this it is not necessary for \( K(X) \) to vanish as \( \dot{\phi} \to 1 \), but only that it be bounded. Other constraints on \( K \) are obtained which allows for a more general framework for the rolling tachyon than was previously considered. The only constraint on the potential is that \( a^3 V \to 0 \) at late times, where \( a \) is the scale factor. The physical motivation is that the tachyon could then be considered as a possible candidate for dark matter [64]. In section §2.4 we discuss the static spherically symmetric solutions to the equations of motion for the most general scalar field Lagrangian with higher derivatives in flat space which are consistent with hyperbolicity and causality. We find the interesting result that for scalar field solutions that are finite at the origin, causality requires its first derivative to vanish there. Furthermore, even though the total energy is positive, the energy density for such solutions is negative at the origin. A physical motivation for
this study arises from the possibility that such scalar fields could describe the dark matter halos around galaxies [63]. In section §2.2 we set up the problem and review some results concerning the criteria for superluminal propagation and hyperbolicity of the scalar field equations. In the concluding section we discuss the results.

2.2 Preliminaries

In this chapter we will be interested in scalar field theories with a Lagrangian of the general form \( L = L(X, \phi) \). Here, \( X = \frac{1}{2} \partial^\mu \phi \partial_\mu \phi \). The equations of motion for the scalar field are given by

\[
\tilde{G}^{\mu\nu} \nabla_\mu \nabla_\nu \phi = L_\phi - 2XL_{XX}, \quad \text{where} \quad \tilde{G}^{\mu\nu} = L_X \eta^{\mu\nu} + L_{XX} \partial^\mu \phi \partial^\nu \phi.
\]

Throughout this chapter we will be using the notation

\[
L_X = \frac{\partial L}{\partial X}, \quad L_\phi = \frac{\partial L}{\partial \phi}
\]

and so on. In (2.1), \( G^{\mu\nu} \) plays the role of an effective metric in which the scalar field propagates. For an equation of this type to have a well defined initial value (or Cauchy) problem and to satisfy global hyperbolicity, the following conditions must hold\(^1\) [60, 61, 62]

\[
L_X > 0, \quad L_X + 2XL_{XX} > 0.
\]

If \( u = 0 \) is the characteristic surface and \( n^\mu = \partial^\mu u \), then the speed of propagation of the small disturbances is given by solving

\[
L_X n^2 + L_{XX} (n^\mu \partial_\mu \phi)^2 = 0.
\]

From this, one deduces the maximum speed to be

\[
n^0 = \frac{W_0 \left( \frac{\bar{\mu} \bar{W}}{|\bar{n}|} \right) + \sqrt{1 + W_0^2 - \left( \frac{\bar{\mu} \bar{W}}{|\bar{n}|} \right)^2}}{1 + W_0^2}.
\]

\(^1\)See appendix §A for further discussion of the Cauchy problem and causality in k-essence
where, \( W_\mu = \sqrt{\frac{L_{XX}}{L_X}} \partial_\mu \phi \). The two cases discussed in this chapter are the time-like spatially homogenous and static spherically symmetric ones. The expressions for the propagation speeds in the two cases are respectively,

\[
\frac{n_0^{\vec{n}}}{|\vec{n}|} = \sqrt{\frac{L_X}{L_X + 2XL_{XX}}}, \quad X = \frac{\dot{\phi}^2}{2}
\]
\[
\frac{n_0^{\vec{n}}}{|\vec{n}|} = \sqrt{\frac{L_X + 2XL_{XX}}{L_X}}, \quad X = -\frac{\phi'^2}{2}.
\]

Here, a dot denotes differentiation with respect to \( t \), and a prime denotes differentiation with respect to \( r \). From these it is easy to see that there is superluminal propagation when \( L_{XX} < 0 \). In summary, the conditions of hyperbolicity and no superluminal propagation may be stated as

\[
(2.2) \quad L_X > 0, \quad L_X + 2XL_{XX} > 0, \quad L_{XX} \geq 0.
\]

For future reference we note that in the static spherically symmetric case when there is no superluminal propagation,

\[
(2.3) \quad \sqrt{\frac{L_X}{L_X + 2XL_{XX}}} \geq 1.
\]

When gravity is included, the effective scalar metric now becomes

\[
(2.4) \quad G^{\mu\nu} = L_X g^{\mu\nu} + L_{XX} \partial^\mu \phi \partial^\nu \phi.
\]

We require this metric to be Lorentzian. In particular in order to have a consistent definition of temporal and spatial directions the largest eigenvalue of (2.4) must be positive while the other three must be negative. This can be shown to be true \([61, 62]\) only if the first two conditions in (2.2) are satisfied (with \( X = \frac{1}{2}g^{\mu\nu} \partial_\mu \phi \partial_\nu \phi \) now) while the last one once again avoids \([58, 60]\) superluminal propagation.
2.3 The Rolling Tachyon

Sen [55] has discussed the qualitative dynamics of a tachyon field in the background of an unstable D-brane system and conjectured that the simplest description within an effective field theory framework can be provided by the following Lagrangian: \( L = V(\phi)K(X) \), with \( \phi \) the scalar field, and

\[
(2.5) \quad K(X) = -\sqrt{1 - \dot{\phi}^2}, \quad V(\phi) = V_0 \exp(-\phi).
\]

The cosmology of this model in the FRW background has been studied in [55, 56, 64], and a particularly surprising result is the existence of solutions with exponentially vanishing pressure at late times, but a nonzero energy density. Since there is no compelling reason for the precise forms in equation (2.5), in this section we keep \( K \) and \( V \) arbitrary (apart from the mild assumptions below) and determine from the constraints (2.2), the conditions under which the equation of state for tachyonic matter becomes \( \omega = 0 \) at late times up to small corrections. The tachyon could then be considered a dark matter candidate in a wider class of models than originally envisioned.

Consistent with the fact that we are dealing with the case of a rolling tachyon, we will make the following assumptions about the potential \( V \) and the kinetic term \( K \): first is that \( K \leq 0 \); second, the range of \( \dot{\phi} \) is bounded; third, \( K \) is bounded as \( \dot{\phi} \to 1^2 \); and fourth, the potential \( V(\phi) \) is positive, has a maximum at \( \phi = 0 \) and monotonically decreases to zero at \( \phi = \infty \) where it is a minimum.

The equations of motion for the scalar field and the scale factor \( a(t) \) are\(^3\)

\[
\ddot{\phi} = -3H \frac{L_X}{L_X + 2XL_{XX}} \dot{\phi} - \frac{\partial \rho_t}{\partial \phi}, \quad H^2 = \rho = \rho_t + \rho_m.
\]

\(^2\)We will take the upper limit of \( \dot{\phi} \) to be 1 in the appropriate units.

\(^3\)We work in units where \( \frac{8\pi G}{3} = 1 \)
For the homogenous FRW background $2X = \dot{\phi}^2 > 0$. $\rho_t = 2XL_X - L$ is the tachyon energy density and $\rho_m$ is that of the rest of matter and $\rho$ is the total energy density. $H = \frac{\dot{a}}{a}$ is the Hubble factor. Note that from the first inequality in equation (2.2) and $L < 0$, it is easy to see that $\rho_t > 0$. Thus, the non-vanishing of the energy density at all times including late times is strictly a consequence of $L_X > 0$. The equation of state parameter for the tachyon is,

\begin{equation}
\omega_t = \frac{L}{2XL_X - L} = \frac{K}{2XK_X - K}.
\end{equation}

Using the $\phi$ field equation of motion it is straightforward to show that

\begin{equation}
\dot{\rho}_t = \frac{d}{dt}(2XL_X - L) = -6HXL_X = -3H(1 + \omega_t)\rho_t.
\end{equation}

The inequality in equation (2.2) then implies that the tachyon energy density is a monotonically decreasing function of time and $\omega_t > -1$. The Hubble factor itself is monotonically decreasing in time as can be seen from

$$
\dot{H} = -\frac{3}{2} ((1 + \omega_m)\rho_m + (1 + \omega_t)\rho_t),
$$

Defining $y = \sqrt{2X}$, and using the factorized form for the tachyon Lagrangian, the tachyon equation of motion may be written as,

\begin{equation}
\frac{dy}{dt} = -\frac{yK_y - K}{K_{yy}} \left(3H \left[\frac{K_y}{yK_y - K}\right] + \frac{\partial V}{\partial \phi} V\right).
\end{equation}

The constraints given in equation (2.2) for the initial value problem to be well defined and the absence of superluminal propagation are expressed in terms of the new variable as:

\begin{equation}
K_y > 0, \quad K_{yy} > 0, \quad K_{yy} > \frac{K_y}{y}.
\end{equation}

Note that whenever $V \to 0$, then $K_y \to \infty$ such that $L_X > 0$. As we will see below, it is this simple fact that guarantees that the tachyon energy density is nonzero and
positive in the limit \( t \to \infty \), while \( \omega \) vanishes; indicating that the tachyon field acts like a non-relativistic dust at late times. Let us consider equation (2.8) at late times. We first discuss the conditions on the potential under which the second term in the brackets is dominant. Let us define the term inside the square brackets in this equation as \( g \), then

\[
\frac{dg}{dy} = \frac{-KK_{yy}}{(yK_y - K)^2}.
\]

Since \( K < 0 \), we see from (2.9) that \( \frac{dg}{dy} > 0 \). The maximum value of \( g \) is thus at \( y = 1 \), which is \( g_{max} \leq 1 \). Moreover, \( H \) is monotonically decreasing. Let us now write for late times \( \phi = t + \theta(t) \), with \( \theta(t) \ll t \). It is easy to check from the above results that the second term inside the brackets in equation (2.8) dominates over the first for late times as long as \( V \to 0 \) faster than \( \frac{1}{a^3} \) when \( t \to \infty \). This condition on the potential will reappear below. Since the overall factor outside the brackets in (2.8) is negative, and since \( \frac{\partial V}{\partial \phi} < 0 \) from our assumptions, the above discussion shows that \( y \) is monotonically increasing as it goes to 1 at late times. In addition, since \( y \) is bounded at \( y = 1 \), then \( \dot{y} = 0 \) at \( y = 1 \). Therefore we conclude that as \( y \to 1 \),

(2.10) \[
\frac{K_y}{K_{yy}} = 0, \quad \frac{K}{K_{yy}} = 0.
\]

As mentioned earlier, \( K_y \to \infty \) for late times while \( K \) is bounded. Thus, the second of the above conditions is not a new requirement since the first implies that as \( y \to 1 \), \( K_{yy} > K_y \). It should be noted that the condition for the absence of superluminal propagation only implies that

\[
\frac{K_y}{K_{yy}} < y,
\]

Thus the condition (2.10) is much stronger.
Let us now expand equation (2.8) about the point $y = 1$ by writing, $\phi = t + \theta(t)$ with $\dot{\theta} < 0$ and $\theta \ll t$. Using (2.10), it is straightforward to get,

$$
\ddot{\theta} = -\dot{\theta}\lambda \left[ 3H + \frac{\partial V}{\partial \phi} \right],
$$

(2.11)

$$
\lambda = \left( 1 - \frac{K_{yy}K_y}{K_{yy}^2} \right) (y = 1).
$$

Integrating this we get (taking $\dot{\phi} \approx 1$ to leading order)

$$
\dot{\theta} = -\alpha a^{-3\lambda} V^{-\lambda},
$$

where $\alpha$ is a constant. Consistency with the boundary conditions require $\lambda < 0$. Since $\theta \ll t$, we see that with a negative $\lambda$, $\dot{\theta}$ vanishes like $a^3 V \to 0$ as $t \to \infty$, which is exactly the condition derived earlier for the term involving the potential $V$ to dominate over the first one in equation (2.8). Thus, $\lambda < 0$, or equivalently, $K_{yyy} > K_{yy}$ at $y = 1$, which is a new constraint on the allowed forms of $K$.

We are now in a position to prove that the equation of state parameter vanishes at $y = 1$ up to small corrections. From equation (2.6) we can obtain,

$$
\frac{d\omega_t}{dy} = \frac{yK_y^2 - K K_y - yK K_{yy}}{(yK_y - K)^2}.
$$

Since $K \leq 0$, and $K_y$ and $K_{yy}$ are both positive, we see that $\omega_t$ is a monotonically increasing function of $y$. Its maximum is therefore at $y = 1$. Near $y = 1$ we can write

$$
\omega_t \approx \frac{K(1) + \dot{\theta}K_y(1)}{K_y(1)}.
$$

However, we have argued above that $K_y(1)$ is infinite, thus $\omega_t = 0$ apart from corrections which vanish like $a^3 V$ at late times.

### 2.4 Background Static Solutions Consistent With Causality

We next consider the static spherically symmetric solutions to the equations of motion of the scalar field in flat space-time. The equation of motion for a static
scalar field $\phi(r)$ is

$$
(2.12) \quad \phi'' + \frac{2}{r} \left( \frac{L_X}{L_X + 2XL_{XX}} \right) \phi' + \frac{L\phi - 2XL_X\phi}{L_X + 2XL_{XX}} = 0.
$$

In the above, $X = -\frac{\phi'^2}{2}$. From section §2.2, the combined constraints of hyperbolicity and absence of superluminal propagation now give the following bound for the coefficient of the $\phi'$ term for all $r$ (see equation (2.3)):

$$
(2.13) \quad \delta = \frac{L_X}{L_X + 2XL_{XX}} \geq 1.
$$

Here we consider only solutions to (2.12) which are finite at the origin. We first want to determine the appropriate boundary condition for $\phi'$ at $r = 0$. We will use a series expansion method for $\phi$ near $r = 0$ to guide us to the correct choice. Even though the coefficient $\delta$ is not a constant but dependent on $\phi$, we know that independent of $r$, $\delta \geq 1$. Therefore, since we are interested in finding the indicial equation in order to determine the boundary condition for $\phi'$, we may treat $\delta$ as a constant. The same applies for the last term in (2.12) as long as we restrict ourselves to solutions which are finite at the origin. These two complications do not affect the indicial equation.

With this in mind, let us look for a series solution of the form

$$
(2.14) \quad \phi = r^s(c_0 + c_1r + c_2r^2 + c_3r^3 + \cdots).
$$

From this we get the indicial equation, $s(s - 1 + 2\delta) = 0$. Since $\delta \geq 1$ for all $r$, we see from (2.14) that in order for $\phi$ to be finite at the origin only the solution with $s = 0$ is allowed. Substituting this expansion into (2.12) we see by matching equal powers of $r$ that $c_1 = 0$. Thus, the boundary condition for this problem that is consistent with (2.13), and the finiteness of $\phi$ at the origin is $\phi' = 0$ at $r = 0$. We now consider the analog of equation (2.7). Let us define $\gamma = -2XL_X + L$. Then
using the equation of motion we obtain

\begin{equation}
\frac{d\gamma}{dr} = -\frac{2}{r} \phi'^2 L_X. \tag{2.15}
\end{equation}

Since $L_X > 0$, we see that $\gamma$ is a monotonically decreasing function of $r$. Therefore, the minimum of $\gamma$ is at infinity. As $r \to \infty$, the solutions to the equations of motion must be such that $\gamma \to 0$ faster than $\frac{1}{r^3}$ in order to keep the total energy content finite. This implies that $\gamma > 0$ as $r \to 0$. From the boundary condition on $\phi'$ at $r = 0$, we see that $\gamma = L > 0$. On the other hand, in the static limit the Hamiltonian density $\mathcal{H}$ becomes $\mathcal{H} = -L$. Thus, we conclude that at $r = 0$, the energy density $\mathcal{H}$ is negative. It is easy to see that the total energy in the static limit is, however, always positive:

$$E = -4\pi \int r^2 L dr = -4\pi \int (\gamma - \phi'^2 L_X) r^2 dr.$$ 

Consider the integral over $\gamma$. Integrating by parts and using the fact that $\gamma$ vanishes faster than $\frac{1}{r^3}$ as $r$ approaches infinity, we get

$$\int \gamma r^2 dr = -\frac{1}{3} \int \frac{d\gamma}{dr} r^3 dr = \frac{2}{3} \int \phi'^2 L_X r^2 dr,$$

where the last equality follows from (2.15). Combining everything, we find that

$$E = \frac{4\pi}{3} \int \phi'^2 L_X r^2 dr,$$

which is manifestly positive definite.

When such a theory is coupled to the Schwarzschild metric, we can look for solutions to the combined equations for both gravity and scalar matter. Such a situation could be relevant for understanding the formation of dark matter halos around galaxies [63]. Though the above analysis has been performed in flat spacetime, our considerations indicate that at least solutions of the scalar field equations
that are finite at the origin should not be relevant to such a scenario if one takes the view that negative energy densities are not allowed. The detailed question of the solutions of the scalar field in the presence of gravity is the subject of the next chapter. Nevertheless, it is interesting that the model we have considered in this section has solutions which violate the weak energy condition at the origin.

2.5 Conclusion

Using the requirement that the field equations are always hyperbolic (and hence the Cauchy problem is well defined) we have obtained a set of consequences for two different problems of physical interest.

For the case of the rolling tachyon in a homogenous FRW background, we have obtained constraints on the forms of the potential and the kinetic terms such that the equation of state of the tachyon vanishes at late times up to small corrections. The tachyon could then be considered a dark matter candidate. The key observation here was that what is required for this to happen is that \( K \) remains bounded but not necessarily zero at late times, while \( K_y \) goes to infinity. The latter in fact guarantees that the energy is non-vanishing in this limit. Other requirements are given by equations (2.10), the potential \( V \) is such that \( a^3V \to 0 \) at late times, and that \( \lambda \) defined in (2.11) be negative. It is easy to check that the choice (2.5) does in fact satisfy all the requirements, but is not unique. The class of models is thus larger than the original.

We have also looked quite generally at the problem of the static spherically symmetric solutions to the equations of motion of the scalar field described by the Lagrangian of section §2.2 and found that if we require the finiteness of the scalar field at the origin, then the solutions consistent with causality have the property that the
energy density becomes negative at the origin. This example brings out very clearly the role that causality plays in the choice of boundary conditions. There have been attempts in the literature [63] to use such scalar field models to describe dark matter halos around galaxies. Clearly, solutions which are finite at the origin will not do the job. However, it is interesting to speculate if this negative energy density at the origin is indicative of an attractive force, analogous to the Casimir effect (but of course classical), at the center of galaxies.
3.1 Introduction

Recently there has been a lot of interest in the applications of non-canonical scalar field theories in cosmology. This is due mainly to the fact that the energy-momentum tensor in such field theories has the potential to describe cosmological fluids with negative pressure, which is a necessary ingredient for accelerated expansion. Examples are: \textit{k-essence} [19], which attempts to explain the accelerated expansion of the universe as well as the coincidence problem, the \textit{DBI action} [27, 28, 57, 35], \textit{tachyonic matter} [55, 56, 1], the \textit{ghost condensate model} [65] and the \textit{Chaplygin gas model} [66]. They are also interesting in the context of inflationary models [32, 3]. Since these theories contain non-standard kinetic energy terms, the constraints of global hyperbolicity and absence of superluminal propagation play an important role [60, 61, 62, 63, 31].

In this chapter we consider static, spherically symmetric solutions of the Einstein equations coupled to a general \textit{k-essence} scalar field. The physical motivation is to look for consistent solutions describing galactic halos [67] and black holes. The standard assumption of galactic dark matter halos is that they consist of an incoherent collection of weakly interacting massive particles. There have been attempts
[63, 68, 69] where some very special type of scalar theories were used to discuss the possibility that the galactic halos could in fact be considered a coherent excitation of a scalar field, much like a boson star. This scenario would have the advantage of providing a unified treatment of dark matter and dark energy since the latter can be described by a k-essence-like theory to begin with.

We consider the most general scalar field Lagrangian that depends on the field φ and invariants of the kind \( X = \frac{1}{2} g^{\mu\nu} \partial_\mu \phi \partial_\nu \phi \), formed from its first derivative. We do not assume that the kinetic terms are separable, nor that they are a quadratic form in the first derivative. Consistent with our aim of finding solutions that describe galactic halos, we impose the generalized “no-force” condition at the origin: i.e., \( \frac{dp}{dr} = 0 \) at \( r = 0 \), where \( p \) is the pressure; and at large \( r \) we demand that the solutions match on to the cosmological solution. We find that solutions do exist that can describe galactic halos, and for certain choices of the metric function, give a good description of the observed rotation curves.

There can be two classes of solutions: those that have negative energy density near the origin and those that don’t. These can be phenomenologically distinguished by the shape of the rotation curves near the origin. Strictly speaking, only one of these solutions has a valid classical description, whereas the other will depend on new physics at short distance scales. Thus, one of the main results of this chapter is that all classical solutions of this theory that satisfy the above mentioned conditions at the origin and asymptotically, must have negative energy density at the origin. This result does not depend on a specific choice of the metric function, but only that it satisfies the conditions for flat rotation curves discussed in [70]. Previous works, for example, the first and last of the references in [68], have assumed a specific form of the metric function \( g_{00} \) for intermediate \( r \), and have found special cases of
the result stated above. We would like to emphasize that our only assumptions are the boundary conditions stated earlier and the general conditions of [70] (see section §3.2). We do not assume anything about the form of $g_{00}$. For the case when the scalar energy density becomes negative near the origin, the total energy or mass enclosed can still be positive. Recent evidence from string theory indicates that there is no justification for restricting ourselves to potentials that are positive everywhere, as long as the total energy is finite [71]. Negative energy densities are not unphysical by themselves, as the Scharnhorst effect [72] clearly shows. Here one has faster-than-light propagation in a Casimir vacuum. In the literature there are studies [73, 74] associating negative energy densities with superluminal propagation of signals. However, in [31] it was argued that superluminal propagation does not necessarily imply the violation of causality for which one requires closed time-like curves. All this clearly deserves further study to see if superluminal propagation without causality violation is consistent. If so, then k-essence-like theories would provide a larger context in which to study the formation of galactic halos.

Another question discussed in this chapter is the existence of black hole solutions in the combined gravity/k-essence system. This problem is of physical interest since galaxies are thought to have black holes at their center. Hence, we would like to ascertain whether black holes can coexist with the halos. Two possibilities are usually mentioned in connection with the issue of black hole-halo coexistence [75]. The first is that the massive black holes were formed together with the galaxies through some internal dynamical process, or secondly, that the black holes are primordial. In the second possibility, they were present even before any luminous activity, and in fact are the source driving the quasars. In both of these situations, once formed, the black holes continue to grow in time. If there is a scalar field pervading the universe
then it can interact with the black holes and if the scalar no-hair theorems [76, 77] for static spherically symmetric black holes are valid it could cause their accretion. One could study this question by considering the stability of the halo in the presence of a background black hole solution. In this chapter, however, we do not do this. Instead our more modest approach is to revisit the scalar no-hair theorems. These are sometimes used to argue [75] that the black holes become heavier in time by “eating” the scalar hair. We present new ways to understand the theorems and suggest some avenues on how they may be circumvented in the context of k-essence like theories. In particular, we clarify the roles played by the choice of asymptotic boundary conditions and the signs of the energy density and the pressure in the proofs [76, 77] of the scalar no-hair theorems.

The chapter is organized as follows. In section §3.2 we discuss the various steps that are necessary to model a spherically symmetric scalar halo and to describe the rotation curves. In section §3.3 we address the question of the no-hair theorems for black holes. We conclude in section §3.4 with a discussion of our results.

3.2 Scalar Fields and Dark Matter Halos

In this section we will describe a halo assuming that it is made up only of scalar fields. Including an exponential disc of baryonic matter should not change the essential conclusions. Thus, we are interested in a scalar field theory coupled to gravity and described by an action of the general form

$$S = \int d^4x \sqrt{-g}L(X,\phi), \quad \text{where} \quad X = \frac{1}{2}g^{\mu\nu}\nabla_\mu\phi\nabla_\nu\phi.$$ 

Its energy-momentum tensor is given by

$$T_{\mu\nu} = \frac{2}{\sqrt{-g}}\frac{\delta S}{\delta g_{\mu\nu}} = L_X\nabla_\mu\phi\nabla_\nu\phi - g_{\mu\nu}L.$$
Here and elsewhere in this chapter, $L_X$ denotes the partial derivative with respect to $X$. The equation of motion for the scalar field is given by

$$\tilde{G}^{\mu\nu} \nabla_\mu \nabla_\nu \phi = L_\phi - 2XL_X\phi,$$

$$(3.1)$$

$$\tilde{G}^{\mu\nu} = L_X g^{\mu\nu} + L_{XX} \nabla^\mu \phi \nabla^\nu \phi.$$ 

The quantities satisfy several constraints which we now outline. An examination of the characteristics of the scalar equation of motion gives the speed of propagation of the scalar fluctuations, and demanding that this so called sound speed ($c_s$) is not superluminal imposes the constraint $\frac{L_{XX}}{L_X} \geq 0$. Demanding in addition that the initial value problem be well posed, and the scalar field equation of motion be globally hyperbolic gives the following list of constraints for this system [60, 61, 62, 63, 31]:

$$L_X > 0, \quad L_{XX} > 0, \quad L_X + 2XL_{XX} > 0.$$ 

In addition, stability requires [31] that $c_s > 0$. Interested readers can refer to appendix §A for further discussion about the origin of these constraints.

As discussed in the introduction, we will be interested in static spherically symmetric solutions of the combined Einstein/scalar system of equations. These solutions must match on to the cosmological solution at large $r$. Depending on the model under consideration, the cosmological solution could either be almost asymptotically flat or not. We will mostly assume almost asymptotically flat solutions, but in the next section we will also consider other possibilities as well. The space-time metric will be described by

$$ds^2 = e^\nu dt^2 - e^\lambda dr^2 - r^2(d\theta^2 + \sin^2 \theta d\phi^2),$$

where $\nu$ and $\lambda$ are functions of $r$ alone. The $(tt)$, $(rr)$ and $(\theta\theta)$ components of the
Einstein equations $G^\mu_\nu = \kappa T^\mu_\nu$ are given by

\begin{align}
(3.2) \quad & \frac{e^{-\lambda}}{r^2}(-1 + e^\lambda + r\lambda') = -\kappa L, \\
(3.3) \quad & \frac{e^{-\lambda}}{r^2}(-1 + e^\lambda - r\nu') = -\kappa(L - 2XLX), \\
(3.4) \quad & -\frac{e^{-\lambda}}{2r}(r\nu'' + \frac{1}{2}r(\nu')^2 + \nu' - \lambda' - \frac{1}{2}r\nu'\lambda') = -\kappa L.
\end{align}

The equation of motion for the scalar field (3.1) in the static limit may be written as

\begin{align}
(3.5) \quad \psi'' + \left[ \frac{\lambda'}{2} + \left( \frac{\nu'}{2} + \frac{2}{r} \right) \frac{L_X}{L_X + 2XLXX} \right] \psi' + \frac{L_\phi - 2XLX\phi}{L_X + 2XLXX} = 0,
\end{align}

where $\psi' = e^{-\lambda}\phi'$, and the speed of sound $c_s$ is given by

$$
\delta(r) = \frac{1}{c_s^2} = \frac{L_X}{L_X + 2XLXX}.
$$

Note that if we require the absence of superluminal propagation of the scalar field fluctuations, then for all $r$, $\delta(r) \geq 1$. Furthermore, stability necessitates that $c_s > 0$.

We are now ready to discuss the question of galactic halos. In [70], it is shown how some essential features of the metric function $\nu$ can be deduced directly from the observed galactic rotation curves, independent of the matter content and a specific gravitational Lagrangian. From stability considerations, [70] concludes that

\begin{align}
(3.6) \quad 0 < r\nu'/2 < 1.
\end{align}

Further assuming circular halos and that information travels to us along null geodesics, one gets an additional constraint [70], which for our purposes may best be written as

$$
|v_c'(r)| < \frac{v_c(1 - v_c^2)}{r},
$$

$$
\frac{r\nu'}{2} = v_c^2.
$$
In the above, \( v_c \) is essentially the tangential component of the velocity [78]; i.e., the rotation curve. These inequalities indicate a smooth scalar profile. It is clear from these that it is not at all difficult to find functions \( v_c(r) \), which give a realistic representation of the rotation curves: one that grows for small \( r \) and flattens out in an intermediate region before matching on to the cosmological solution outside. This is precisely the kind of behavior we demand of the corresponding derivative of the metric function \( \nu \), so that our model is able to describe the observed rotation curves.

We now turn to the main part of this section; namely, a discussion of the halo-like solutions of the Einstein/scalar system subject to the appropriate boundary conditions at \( r = 0 \), and at \( r \to \infty \). Asymptotically, the solution must match on to the cosmological solution. However, in the absence of an explicit cosmological solution we will assume that the pressure goes (almost) to zero as \( r \to \infty \); i.e., the condition of asymptotic flatness. We would next like to determine the appropriate boundary conditions for small \( r \). Since we are describing a halo as a coherent state of a scalar field, the appropriate boundary condition at the origin is the generalized “no-force” condition; i.e., \( \frac{dp}{dr} = 0 \), where \( p \) is the pressure and \( \frac{dp}{dr} \) is given by the Oppenheimer-Volkov equation. For our model, the pressure and the energy density are \( p = L - 2XL_X \) and \( \rho = -L \), respectively. Using (3.5) or (3.8), it is straightforward to obtain \( \eta' = e^{-\frac{1}{2}\phi'} \)

\[
(3.7) \quad \frac{dp}{dr} = -\frac{1}{r}(\frac{r\nu'}{2} + 2)\eta'^2L_X.
\]

Since the quantity inside the brackets is bounded between 2 and 3, we obtain the boundary condition for the scalar field equation, which is that \( \eta' \) goes to zero faster than \( \sqrt{r} \) at the origin (keep in mind that \( L_X > 0 \) everywhere).

In order to understand the type of solutions one can obtain, we will adopt the
following strategy: we will first look for consistent solutions at the origin and follow these to \( r \to \infty \) using (3.7). We will not use any specific form for the metric functions other than the fact that for halo-like solutions the component \( g_{00} \) satisfies (3.6) for all \( r \). Our conclusion will be that under these conditions, the only non-trivial solutions consistent with the boundary conditions are those that have negative energy density at the origin. Thus, our aim is not to find explicit solutions but to show this universal feature of all halo-like solutions under a very general set of assumptions.

In accordance with this program, let us next look for consistent solutions of the different equations of motion near the origin. First, consider the scalar field equation (3.5). Using the variable \( \eta' \) \( (e^{-\lambda} \neq 0 \) near the origin), it may be rewritten without explicit reference to the metric function \( \lambda \):

\[
\eta'' + \frac{1}{r^2} (\frac{r \eta'}{2} + 2) \delta(r) \eta' + h(r) = 0,
\]

(3.8)

where \( h(r) = \frac{L_{\eta} - 2X L_{X\eta}}{L_X + 2X L_{XX}} \).

Consider a small \( r \) expansion for \( \eta \) and write

\[
\eta = r^s(c_0 + c_1 r + c_2 r^2 + \cdots).
\]

(3.9)

As we will see next, a consequence of (3.8) and the condition on the vanishing of \( \eta' \) at the origin is that there are two possibilities for \( h(r) \): (1) \( h(r) \) is smooth at the origin; or (2) it behaves as \( \frac{1}{r^\epsilon} \), where \( 0 < \epsilon < 1/2 \). To see possibility (1), we will first assume the opposite; i.e., suppose that \( h(r) \) goes like \( \frac{1}{r^\epsilon} \) to leading order at small \( r \). The equation of motion (3.8) implies that to leading order, \( \eta' \) approaches a nonzero constant at small \( r \), which rules out this possibility. Similarly, if \( h \sim 1/r^n \), where \( n \geq 2 \) is an integer, then \( \eta' \sim 1/r^{n-1} \), which would contradict our previous finding that \( \eta' \to 0 \) at the origin. As we discuss below, in this case \( \eta \to constant \) at small \( r \). The second possibility arises when \( \eta \sim r^b \) with \( 3/2 < b < 2 \) for small \( r \).
The condition that $b > 3/2$ follows from the requirement that $\eta'$ vanish at the origin faster than $\sqrt{r}$. To see that $b < 2$, we have to use the fact that at the origin $p = L$, and therefore $\frac{dL}{dr} = 0$ there. Suppose $L \sim \eta^a$ ($a > 0$, since $a = 0$ is included in case (1)), then the no-force condition implies that $ba > 1$. Since both $L_X$ and $L_{XX}$ are positive definite, $h(r) \sim L_\eta \sim r^{ba-b}$. Now, from the equation of motion (3.8), we see that $ba - b = b - 2$. We need consider only the case $a < 1$, since the other possibility is included in case (1). Using $a = \frac{2b-2}{b}$, it is easily seen that $b < 2$. We will consider each of the two cases separately starting with case (1).

For case (1), the only allowed behavior of $h(r)$ consistent with the boundary conditions is that it is well defined at the origin. Substituting (3.9) into (3.8), we find for the indicial equation ($c_0 \neq 0$)

$$s(s - 1 + \bar{\gamma}\bar{\delta}) = 0,$$

where $2 < \bar{\gamma} < 3$.

In the above, $\bar{\gamma}$ denotes the value of $\gamma = rv'/2 + 2$ for small $r$, and the inequality above follows from (3.6). Further, if we demand the absence of superluminal propagation, then $\bar{\delta}$ (the value of $\delta$ at small $r$) is greater than 1. Thus, we see that the absence of superluminal propagation forces upon us the solution with $s = 0$. The actual behavior of $\eta'$ near the origin depends on the small $r$ behavior of $h(r)^1$. By substituting the expansion (3.9) with $s = 0$ into (3.8), we find that $c_1 = 0$. Therefore, the solutions at small $r$ behave like $\eta \sim \text{constant}$, and $\eta' \sim r$. We will see below how these conditions translate into the shape of the rotation curves at small $r$.

We must now find a consistent solution for small $r$ for the metric functions as well. For these we turn to the Einstein equations (3.2) and (3.3). It is useful to rewrite

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1To consider an extreme case, if the Lagrangian depends only on $X$, then this term is absent. It is then trivial to verify that the only solution which is smooth at the origin is $\eta = \text{constant}$. This in turn implies the constancy of the pressure and the energy density in this case. This is the de Sitter solution. We will not discuss this possibility further.
these as

\begin{align}
\lambda' &= -\kappa rL - (e^\lambda - 1)(\kappa rL + \frac{1}{r}), \\
\nu' &= \kappa rp + (e^\lambda - 1)(\kappa rp + \frac{1}{r}).
\end{align}

On the right hand side of these equations we have the energy density $\rho = -L$ and the pressure $p = L - 2XLX$. From these we see that the combination $\nu' + \lambda'$ only depends on $p - L$, which is proportional to $X$ (see also (3.17) below). The functions $\nu$ and $\lambda$ themselves depend on $p$ and $L$ separately. In general, the small $r$ behaviors of $L$ and $(p - L) \sim X$ can be different: $L$ is a function of both $X$ and $\phi$, and $\phi$ goes to a constant at small $r$. From a Taylor expansion of $L$ around $X = 0^2$, it is seen that the leading order behavior of $L$ can therefore be either a constant, or that of at least $X$. From (3.10) and (3.11) we see that in the former case, the leading order behavior of $\nu'$ and $\lambda'$ are determined from that of $L$, and this leading order behavior is cancelled from the sum. It is the subleading terms of $\nu'$ and $\lambda'$ that are now proportional to $X$. For small $r$, we have two possibilities: (i) both $L$ and $2XLX$ have similar behavior near the origin; i.e., since $L_X > 0$, then $L \sim 2XLX \sim r^2$ at most; or (ii) $L \sim constant$, and $2XLX \sim r^2$ or faster. In the case of possibility (i), the scalar field and in particular the speed of sound plays an important role in restricting the small $r$ behavior. In contrast, for possibility (ii) the leading order behavior is governed by $L$, and the constraints on the absence of superluminal propagation do not play a role in determining the small $r$ behavior of the metric functions. We will discuss each possibility below.

For possibility (i), it is straightforward to check that consistency with the Einstein equations gives the following behaviors at small $r$ for the metric functions: $\nu, \lambda \sim r^4$ or faster. At this point the scalar energy density and pressure go like $r^2$ or faster.

\footnote{We assume that the Lagrangian is Taylor expandable in $X$.}
This behavior excludes a purely $X$ independent potential term in the Lagrangian for the scalar field, since that could lead to a constant behavior for $\rho = -L$ at small $r$.

Let us see how the solutions we have discussed at $r = 0$ match on to the solutions at $r \to \infty$. Equation (3.7) is useful for this purpose as well. The term in brackets on the right hand side of this equation is positive, so $p$ is monotonically decreasing as we go from some small $r$ to infinity, where it approaches the (almost zero) value dictated by the near asymptotically flat cosmological solution outside the galactic halo. The pressure at $r = 0$ is zero since, as discussed above, in this case $p = L - 2XLX \sim r^2$. Since $p$ is monotonically decreasing to zero for all $r$, one should conclude that the only possibility is the trivial solution, $p = 0$ everywhere. While this would be true if the classical theory is valid everywhere, it is possible that quantum gravity effects modify the theory at small scales, negating this argument. The influence of the requirement of subluminal propagation has in fact a very indirect influence on all this. The real reason for the absence of the negative energy density is that demanding the leading small $r$ behavior of $\eta'$ to be relevant both for the pressure and the Lagrangian, eliminates the pure potential terms from consideration. We will discuss this in more detail below.

 Possibility (ii) leads to a completely different conclusion. In this case it is easy to check that the metric functions have the following behavior for small $r$: $\nu, \lambda \sim r^2$. Let us now see what happens as $r \to \infty$. As was the case in possibility (i), equation (3.7) implies that the maximum positive value of $p$ is for small $r$. The asymptotic value of $p$ is zero, therefore, $p \approx L > 0$ at the origin. Thus, the energy density $\rho = -L$ can now be negative for small $r$. The Lagrangian can contain terms that depend only upon $\phi$ and not its derivatives. Such potential terms can approach a (negative) constant at small $r$. This case is analogous to the flat space situation discussed in [1].
we will see below, this should not be surprising since the condition for the absence of superluminal propagation has not played a role here. Even though the energy density can be negative in some region, we will now see that under very reasonable assumptions, the total mass within a large enough region is positive definite. The metric function \( \lambda \) is related to the mass function \( m(r) \) by

\[
e^\lambda = (1 - \frac{2Gm}{r})^{-1}.
\]

Consider the total mass inside a radius \( R \). From (3.2) and (3.12), we have

\[
m = -4\pi \int_0^R L r^2 dr = -4\pi \int_0^R (p + 2XL_X)r^2 dr.
\]

For the first term in the integral we perform an integration by parts to write

\[
\int_0^R pr^2 dr = -\frac{1}{3} \int_0^R \frac{dp}{dr} r^3 dr + \frac{1}{3} [pr^3]_0^R.
\]

Using (3.7), and substituting into (3.13), we obtain for the mass parameter \( m(R) \)

\[
m = \frac{4\pi}{3} \int_0^R (\eta')^2 L_X (1 - \frac{ru'}{2}) r^2 dr - \frac{4\pi}{3} [pr^3]_0^R.
\]

Using (3.6) the first term is manifestly positive definite, and the surface term can be made small for large enough \( R \) if the pressure falls off faster than \( \frac{1}{R^2} \) (asymptotically near flat condition), or if the cosmological solution is such that for large \( r \) the pressure is negative, as is the case with some k-essence models.

Let us now take up case (2); i.e., \( h(r) \sim \frac{1}{r^\epsilon} \), where \( 0 < \epsilon < 1/2 \). As we have discussed earlier, in this case \( \eta \sim r^b \) with \( b > 3/2 \) and \( L \sim \eta^a \) with \( 0 < a < 1 \). Thus, at the origin, \( p = L \to 0 \). The situation here is similar to the one discussed for possibility (i) above. When we try to match this behavior to that of the solution at infinity using equation (3.7), we will again conclude that the only consistent solution is \( p = 0 \) everywhere.
We have therefore arrived at the conclusion we had mentioned earlier. The only non-trivial halo-like solutions using the classical equations of motion of the Einstein/k-essence system with the stated boundary conditions are those with negative energy density at the origin.

The origin of the negative energy densities in certain regions can be clearly traced to the pure potential terms in the Lagrangian. Indeed, as $\phi$ goes to a constant for small $r$, these potentials tend to a negative constant. If the total mass parameter is positive, then it appears to us that excluding theories where the potential can be negative in some regions is not justified. In fact, potentials such as these have been recently considered in a variety of situations. For example, in supersymmetric AdS compactifications one encounters potentials with local negative maximums. Additionally, a large class of supersymmetric compactifications, including Calabi-Yau and G2, give rise to effective four dimensional potentials with negative regions [71]. The stability of the solutions considered in this chapter will be the subject of a future investigation. We would like to note, however, that there are many known examples where the potential is negative at an extremum and yet the solution is stabilized due to gravitational effects [79] as long as the scalar field theory satisfies the positive energy theorem [80].

Another reason to exclude negative energy densities is its association with superluminal propagation and causality violation. We now turn to a discussion of this point. Consider the constraints on superluminal propagation in some more detail. Until now the only constraint we have considered is the one arising for the speed of the scalar fluctuations, which is obtained from an analysis of the characteristics of the scalar field equation of motion. A simple argument shows that this has no direct connection with the sign of the energy density. If one changes $L \rightarrow -L$, the
scalar equations of motion do not change, and hence neither does the analysis of the characteristics, but the energy density changes sign. Therefore, there must be an additional constraint on superluminal propagation, which is dependent on the sign of the energy density. In fact, there have been quite a few papers which investigate the relationship between faster-than-light travel and negative energy densities [73]. They all differ on the precise definition of superluminal propagation. For the purposes of this chapter we follow the discussion of [74], which is specific to static spherically symmetric space-times for which the Killing time can be used to measure the time required for objects or signals to propagate between two of its orbits. In [74] a theorem was proven that if in such a space-time, the (time-like) weak energy condition is satisfied, then the signaling time is never faster than the corresponding signal in Minkowski space. The normalization of the Killing time is appropriate for an observer at very large distances. More specifically, the absence of superluminal signals in the sense defined above requires that $T_{\mu\nu}t^\mu t^\nu \geq 0$ for any time-like vector $t^\mu$. Demanding the absence of superluminal propagation based on this criterion would eliminate possibility (ii), since it implies negative energy densities at very small $r$. However, superluminal travel by itself is not threatening as long as there is no causality violation for which closed time-like curves are required. In fact recently [31] it was realized that, from the viewpoint of pure classical field theory, models which allow for superluminal propagation even on dynamical backgrounds do not necessarily possess internal contradictions. In particular, these models do not lead to any additional causal paradoxes over and above those already present in standard general relativity.

We will therefore take the point of view that the solutions discussed above cannot be a priori discarded (if the total mass parameter is positive definite and well be-
haved) and look to the data to decide if both possibilities are realized or not. In this regard it is interesting to note that the two possibilities discussed above can predict very different behaviors for the rotation curves at small \( r \) depending on the model: possibility (ii) implies that \( v_c \sim r \); and if such a possibility exists in the complete theory, (i) implies that \( v_c \sim r^2 \) or faster. Our discussion suggests that the shape of the rotation curves at small \( r \) should therefore provide the phenomenological distinction between the possibilities considered.

If the theories we consider can indeed describe halos of dark matter, then we need to have an understanding of its interactions with a black hole, which presumably are at the center of galaxies. With this in mind, the next section will revisit the “no scalar hair” theorems [76, 77] for black holes.

### 3.3 Dressing Black Holes With Scalar Fields

In this section we study the possibility of a scalar field dressing an asymptotically flat, static, and spherically symmetric black hole solution in the theory discussed in the earlier sections. As discussed in the introduction, this question is of relevance for the stability of the halos described previously.

The distinguishing feature of a black hole is the existence of an event horizon whose position at \( r_s \) is determined from the condition, \( g^{rr}(r_s) = 0 \). Physical quantities, in particular the components of the energy-momentum tensor, are regular at the horizon, from which we deduce the regularity of the scalar field and its derivative at \( r = r_s \), instead of regularity at the origin. With this as the boundary condition, let us consider the scalar field equation of motion near the horizon. We want to find the behavior of the allowed solutions consistent with the absence of superluminal propagation. From equation (3.5) we see that we need to know the behavior of \( \nu' \)
near the horizon. For this purpose let us consider the \((tt)\) component of the Einstein equation (3.2) near \(r = r_s\) and write

\[
e^{-\lambda} = A_\lambda (r - r_s) + \mathcal{O}((r - r_s)^2).
\]  

(3.14)

Thus,

\[
-\lambda' = \frac{1}{r - r_s} + \text{constant} + \mathcal{O}(r - r_s),
\]  

(3.15)

and from (3.2) we obtain

\[
A_\lambda = \frac{1 + \kappa r_s^2 L(r_s)}{r_s} > 0.
\]  

(3.16)

The last inequality follows from the fact that \(e^{-\lambda}\) must be growing as we move out from the horizon. There is the possibility that \(A_\lambda = 0\), which would give the leading order behavior in (3.14) as \((r - r_s)^2\). This would correspond to the case of an extremal black hole solution, which we will not discuss further. We are now able determine the leading order behavior of \(\nu\) near the horizon consistent with the Einstein equations (3.2, 3.3). From these

\[
\nu' + \lambda' = e^\lambda r \kappa f(r),
\]  

(3.17)

\[
f(r) = -2 XL_X > 0.
\]  

(3.18)

Integrating this in the neighborhood of the horizon we get

\[
\nu + \lambda = \frac{r_s}{A_\lambda} f(r_s) \ln(r - r_s) + K + \mathcal{O}(r - r_s),
\]  

where \(K\) is an integration constant. Using the previously determined value of \(\lambda\), we finally obtain the leading order behavior of \(\nu\), assuming that \(A_\lambda \neq 0\):

\[
\nu = B \ln(r - r_s) + \text{constant} + \mathcal{O}(r - r_s), \quad \text{where} \quad B = \frac{1 + \kappa r_s^2 (L(r_s) + f(r_s))}{1 + \kappa r_s^2 L(r_s)}.
\]
It is easy to check from the finiteness of the Ricci scalar at the horizon that $B = 1$, which tells us that $f(r_s)$ approaches zero there. Since $L_X > 0$, we see that $X \to 0$ as $r \to r_s$. The exact rate at which this happens will be determined below. From this, we see that

$$\nu' = \frac{1}{r-r_s} + \text{constant} + \mathcal{O}(r-r_s),$$

$$e^\nu = A \nu(r-r_s) + \mathcal{O}((r-r_s)^2),$$

where the constants in (3.19) and (3.15) are not related and $A_\nu > 0$.

Let us now look for a series solution of the scalar field that is regular at the horizon. In addition, we have seen from the discussion immediately preceding equation (3.19), that the solution should be such that $X = -\frac{e^{-\lambda \phi^2}}{2}$ must vanish there. Introducing $\delta r = r - r_s$, we have the following expansion:

$$(3.20) \quad \psi = \delta r^s(a_0 + a_1 \delta r + a_2 \delta r^2 + \cdots).$$

Substituting this into equation (3.5) we get the following indicial equation:

$$s(2s + \alpha - 3) = 0,$$

where $\alpha = \frac{1}{c^2(r_s)}$. From the discussion on subluminal propagation in the previous section, $\alpha \geq 1$. In this instance, both the solutions for $s$ are not immediately excluded, provided that $\alpha$ is not greater than 3. The special case of $\alpha = 3$ is the same as the $s = 0$ case, except the second solution has a leading log $r$ behavior.

Let us first consider the possible solutions with near horizon behavior dictated by $s = 3s - \alpha/2$. As a result of the smoothness of the action near the horizon, the small $\delta r$ leading order behavior of $\psi'$ is $\delta r^{1/2-\alpha/2}$. This behavior is too singular if we are considering only subluminal propagation, and hence, this solution is ruled out. Next, let us consider the solution with $s = 0$. For the same reasons discussed in the previous
section, \( a_1 = 0 \), and the leading order behavior of \( \psi' \) near the horizon is of the form \( \delta r \). Moreover, this near horizon behavior implies that \( \phi' = \text{constant} \), or equivalently \( X \sim e^{-\lambda} \phi'^2 \sim \delta r \). This is completely consistent with the Einstein equation (3.17) (Recall that always, \( L_X > 0 \)). Indeed, from the combined near horizon behavior of \( \nu' \) and \( \lambda' \) we see that the leading term on the LHS of this equation is a constant, which is the same on the RHS. We should comment that the Einstein equation (3.17) rules out the possibility \( s = 3/2 - \alpha/2 \) when \( \alpha < 1 \) (superluminal propagation). Thus, it should be emphasized that consistency (at the horizon) for the \( s = 0 \) solution holds irrespective of whether we have superluminal or subluminal propagation.

We must now check the consistency of the solution at large \( r \). The main question we would like to address is whether the solutions at the horizon can match on to the ones at infinity in a manner consistent with the Einstein equations. First let us list the behavior of the pressure \( p \) and \( \frac{dp}{dr} \) at the horizon and asymptotically. We will see how far one can go without assuming the weak energy condition. Consider

\[ dp \over dr = -\left(\nu' + \frac{2}{r}\right)\eta'^2 L_X, \tag{3.21} \]

near the horizon. Using (3.19) and (3.14), and the fact that \( \phi' = K \) (a constant), we get

\[ dp \over dr \sim -\frac{1}{2}A_\lambda K^2 L_X < 0. \]

In order to find its behavior at infinity we need the asymptotic form of \( \nu' \). For this purpose, one could first find the asymptotic form of \( \lambda \) from the integrated version of (3.2):

\[ e^{-\lambda} = 1 - \frac{2Gm(r, r_s)}{r} - \frac{r_s}{r}, \quad \text{where} \quad m(r, r_s) = 4\pi \int_{r_s}^{r} \rho r'^2 dr'. \]
Assuming that $\rho$ falls off faster than $\frac{1}{r^3}$ at large $r$, the leading order behavior of $\lambda'$ is

$$-\lambda' \sim \frac{2Gm(r,r_s)}{r^2} + \frac{r_s}{r^2}. $$

From (3.17) and (3.18), it would follow that

$$\frac{\nu'}{2} > \frac{(Gm(r,r_s) + \frac{r_s}{2})}{r^2} \geq 0,$$

which in turn would imply from (3.21) that $\frac{dp}{dr} < 0$ also for large $r$. However, we want to get information about the large $r$ behavior of $\frac{dp}{dr}$ without using the that $\rho$ falls off faster than $\frac{1}{r^3}$ at large $r$. For this purpose we note that since $\frac{dp}{dr} < 0$ near the horizon, in order for it to be positive asymptotically at large $r$, it would have to vanish in between. We will now show that this is untenable on the basis of the Einstein equations, apart from the trivial solution that $\phi$ is a constant. For this purpose note that from (3.21), when $\frac{dp}{dr} = 0$

$$\frac{\nu'}{2} = -\frac{2}{r},$$

It is easily checked using (3.3) that this implies that

$$e^{\lambda} = -\frac{3}{1 + \kappa pr^2}. $$

Since we are looking for regular black hole solutions of the combined Einstein/scalar system, we must rule out the possibility of a change in signature of the metric, which means that $1 + \kappa pr^2 < 0$ or $\kappa pr^2 < -1$. From (3.17) and (3.18), however,

$$\frac{-2 + 3\kappa pr^2}{r(1 + \kappa pr^2)} > 0,$$

or $\kappa pr^2 > -\frac{2}{3}$. The two constraints on $\kappa pr^2$ are incompatible, hence we conclude that $\frac{dp}{dr}$ cannot vanish between the horizon and infinity and must therefore also be negative asymptotically. We emphasize that our conclusion that $\frac{dp}{dr}$ is negative
everywhere outside the horizon does not rely on either the weak energy condition or the asymptotic flatness condition.

Let us next find how the pressure \( p(r) \) behaves near the horizon and asymptotically. To find \( p \) near the horizon, note here that \( \eta^2 \sim (r - r_s) \). Thus, at the horizon, \( p = -\rho \). This is the only place we need to worry about whether \( \rho \) is positive or negative. If \( \rho \) is negative then \( p \) is positive at the horizon, and if instead \( \rho \) is positive then \( p \) is negative there. Let us now consider the possibilities for \( p \) at large \( r \). For positive \( \rho \), since \( \frac{dp}{dr} \) cannot change sign in between, \( p \) must also be negative at large \( r \). Now it is easy to see that if such a solution were to exist, it cannot be asymptotically flat. Indeed, writing \( \eta^2 L_x = p - L \) in (3.21), and integrating using the integrating factor we get

\[
p = -\frac{e^{-\frac{2}{r^2}}}{r^2} \int_{r_s}^{r} (e^{\frac{2}{r^2}}) \rho dr.
\]

Let us see what this implies for \( p(r) \) at asymptotic values of \( r \). Using the asymptotic condition that \( \rho \) falls off faster than \( \frac{1}{r^3} \) at large \( r \), we see that the integral converges and \( |p| \) falls off at least like \( \frac{1}{r^2} \) at large \( r \). However, since we have already argued that \( \frac{dp}{dr} \) is negative, it must be that \( p \) is positive asymptotically at large \( r \). The condition of asymptotic flatness gives a positive value for \( p \) at large \( r \), thus excluding this possibility. Furthermore, asymptotic flatness is incompatible with the constraint \( L_x > 0 \). This is a result of the fact that \( L_x > 0 \) implies that \( \rho > -p \). Since \( p \) is negative in this case and monotonically decreasing, it follows from \( \rho > -p \) that \( \rho \) has to be positive and monotonically increasing in order for \( L_x > 0 \) to be satisfied. The only way we can have an asymptotically flat solution is if \( p = 0 \) everywhere. We therefore conclude that for positive energy density \( \rho \), the only allowed black hole solution has negative pressure asymptotically at large \( r \), and does not obey the asymptotic flatness condition. In k-essence models it is possible to
have a cosmological solution in the matter dominated epoch with these properties. However, it is not clear how we can match on a dynamically evolving scenario onto a static solution.

When the energy density is negative, the pressure can be both positive or negative for large $r$ indicating a solution where one can get either asymptotic flatness or not. However, if we invoke the results of [74] regarding superluminal propagation of signals discussed in the previous section, then this possibility would be excluded. For reasons mentioned earlier, we do not subscribe to this viewpoint as long as the enclosed black hole mass is positive and well behaved. We should mention at this point that solutions that can have negative energy density somewhere should be checked for stability.

Our discussion in this section suggests that the possibility of a black hole and a scalar halo forming at the same time is difficult. This however, does not necessarily make the scalar halo model problematic. It is quite possible, and there are other evidence [75] in support of a scenario where a primordial black hole was present before the formation of the galaxy itself, triggering the quasar activity. The halo can then be treated as a perturbation on this primordial black hole background. An analysis of this is beyond the scope of this chapter and is the subject of an ongoing project.

3.4 Conclusion

We have found some consistent solutions of the gravity/scalar system, which can describe galactic halos. Solutions which have negative energy density near the origin have rotation curves with $v_c \sim r$ at small $r$. Classically these are the only allowed solutions, however, should solutions that have vanishing pressure at the origin be
allowed in a complete theory, the rotation curves should show much steeper behavior for small $r$. Solutions with negative energy density are associated with superluminal propagation as discussed in section §3.2, however, that does not necessarily imply causality violation. The total energy of these configurations can be positive definite even if the energy density is negative somewhere. On the basis of this, and consistent with similar phenomenon in other physical situations [71], we have argued against excluding such configurations a priori. Our analysis in section §3.2 is based on spherical halos. However, under reasonable assumptions, we expect similar conclusions in general. This is because we can parametrize the departure from sphericity by a parameter which should have a smooth limit to zero unless there are topological obstructions.

We have also reconsidered the question of static spherically symmetric black hole solutions in a theory of gravity coupled to scalar fields with non-standard kinetic energy terms. We have considered situations where the scalar energy density can be negative in some regions, but the total mass is still positive, as well as the almost asymptotically flat and other boundary conditions. In the case when the (time-like) weak energy condition is satisfied and the asymptotically flat boundary conditions are enforced, we recover the scalar no-hair theorems [76, 77]. Our analysis is based on the minimum of assumptions and different from previous ones. We find loopholes not only for the case of negative energy densities, but also when the boundary conditions are not asymptotically flat. For reasons mentioned earlier, we do not think it is justified to, a priori, discard solutions where the scalar energy density is negative somewhere but where the total black hole mass is positive and well behaved. In this context it should be mentioned that in [33] consistent black hole solutions have been found in the kind of theories considered here which are stationary and do contain
regions where the energy density can be negative. We have not discussed the stability of any possible solutions. This, and the construction of explicit solutions is also an ongoing project.
CHAPTER IV

Stationary Solutions

4.1 Introduction

The discovery of the late time acceleration of the universe using Supernova Ia [9] confirmed by other observations ([81] and references therein), opened a window of opportunity for the existence of novel cosmological scalar fields not only during the early inflationary stage but also in the current universe. Indeed, scalar fields are the most natural candidates for the realization of inflation, and for the dynamical explanation of dark energy (DE) that is responsible for the late time acceleration. Arguably, the main difficulty in the modeling and understanding of the possible dynamics of dark energy arises because of the fine tuning issues. In particular, there is the so-called coincidence problem [15, 29]: why is DE only now comparable with the energy density in the dust-like dark matter? This coincidence would be especially remarkable, if one assumes that both these dark constituents are independent of each other and evolve very differently in time. Partially because of the fine tuning problems it is not surprising that the candidates for DE often have not only rather exotic names: quintessence/cosmon [15, 16], k-essence [19, 82], phantom [83, 84], ghost condensate [65], quintom [85], etc, but also very unusual properties. In particular, these scalar fields can possess: extremely small effective mass (quintessence,
quintom), sound speed which can be much smaller and even larger than the speed of light (k-essence, ghost condensate), negative kinetic energies (phantom, quintom), Lorentz symmetry breaking and gravity modifications even around the Minkowski space-time background (ghost condensate). The most successful paradigm to solve the coincidence problem is currently k-essence, where the highly nonlinear dynamics triggers the equation of state of DE from radiation-like to quasi de Sitter around the transition to the matter domination stage. In the late matter domination epoch the k-essence has a speed of sound which is much smaller than one. However, it was showed [30, 86] that to explain the coincidence problem k-essence models must necessarily have at least a short phase where the fluctuations in the k-essence travel at superluminal speeds. For our purposes it is important that the nonlinear dynamics responsible for the attractor behavior addressing the coincidence problem requires an explicit dependence of the Lagrangian on the field strength [86]. This field dependence cannot be eliminated by any field redefinitions. Thus, successful k-essence models as well as quintessence/cosmon models cannot be shift symmetric.

On the other hand it is known that the current universe is highly inhomogeneous on small scales and in particular that there are plenty of black holes (BHs) of different mass and origin. Thus an interesting and natural question arises, how do black holes surrounded by cosmological scalar fields evolve? In addition, from the theoretical viewpoint it is interesting to consider BHs “dressed” with different field backgrounds. This could have a valuable impact on our understanding of the physics of horizons (see e.g [33, 87, 88]). Owing to the no-hair theorems [77, 89] we know that BHs cannot support static configurations of scalar fields\(^1\). Therefore, any scalar hair will

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\(^1\)BHs can not support scalar hair at least for theories that respect some of the standard energy conditions. Having in mind the exotic properties of DE models mentioned above, it would be interesting to find examples of stable scalar hair in theories violating the usual energy conditions. For a model of hairy scalar BHs with ghost like quantum instabilities see Ref. [90].
be continuously swallowed by the BH. In particular one could analyze the growth (and may be even formation) of black holes due to the accretion (collapse) of DE. Then one can try to use powerful and rather universal laws of black hole thermodynamics, combined with astrophysical observations to restrict the allowed properties of DE candidates and rule out some of them as contradicting either BH thermodynamics or astrophysical data. Recent studies along these lines were, for example, done in [75, 33, 87, 88, 91, 92, 93, 94, 95].

Finally, for k-essence, a typically very small sound speed during the late matter domination era allows for rather significant large-scale inhomogeneities around BHs and other massive objects. This long-range clumping would be one of the characteristic, potentially observable consequences of k-essence. Moreover, due to this ability to realize small sound speeds along with the dust-like equation of state, the k-essence fields can be used to model dark matter [63, 69, 2]. In this setup, the presence of supermassive BHs at the center of galaxies makes understanding the accretion process even more necessary.

On the other hand the presence of backgrounds with the superluminal sound speed mentioned above opens an exciting possibility to look beyond the BH horizon [33]2. Note that the current bounds [97, 98] on DE sound speed are not restrictive at all.

The classical and most simple setup for accretion problems is a steady state or Bondi accretion [99]. Remarkably, a lot of astrophysical phenomena can be described by a steady state accretion. For a review see e.g. [100]. For scalar fields, the Bondi accretion was recently studied in [33, 93, 94, 95, 101]. It is fair to say that almost all known analytical solutions3 for accreting scalars either belong to the Bondi case or

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2Despite of the presence of superluminal propagation, the accretion backgrounds constructed in these works are free of any causal pathologies [31]. However, it is interesting to study whether, similar to Ref [96], two boosted BH could create causal paradoxes in this setup.

3See however Ref [75]
represent the dust-like free fall. The dust-like time dependent accretion of a massive canonical scalar field was considered in [95], while dust-like solutions for the ghost condensate scalars were found in [92, 88]. It seems that scalar fields with canonical kinetic terms would not leave any important impact on the astrophysical BHs in the current universe [95]. Nevertheless, accreting scalar fields could play an important role in the formation of primordial BHs (see e.g. [75]).

In this chapter we investigate stationary configurations for general k-essence scalar field theories. We show that the necessary condition for the existence of exact stationary configurations is the symmetry of the theory with respect to constant shifts in the field space: \( \phi \rightarrow \phi + c \). This symmetry has to be realized either in terms of the original field strength or after a field redefinition. On the way, we also analyze properties of general k-essence scalar field theories covariant with respect to field redefinitions. The proof is valid for general theories with nonlinear kinetic terms in both the test-field approximation and the self-consistent case where the background metric is governed by the field \( \phi \) itself. It is interesting to note that shift symmetric scalar field theories are exactly equivalent to perfect fluid hydrodynamics provided that only such field configurations which have time-like derivatives are considered. In particular this result implies that the most interesting scalar field models of dark energy cannot realize a steady state/Bondi accretion. Thus, in general, the solution to the problem of accretion of these fields onto black holes requires a knowledge of their initial configuration. In this chapter we are discussing stationary configurations that are exact. Of course, in the real world the stationarity should be considered an approximation. It may well happen that the solutions would only asymptotically approach the stationary regime. For some canonical scalar fields this behavior was demonstrated in [95].
4.2 Derivation of the stationary configurations

Let us consider a general scalar field theory with the action

\[ S = \int d^4x \sqrt{-g} L(X, \phi), \quad \text{where} \quad X = \frac{1}{2} g^{\mu\nu} \nabla_\mu \phi \nabla_\nu \phi. \]

Here, \( g_{\mu\nu} \) is the gravitational metric and as usual \( g \equiv \text{det} g_{\mu\nu} \). Throughout the chapter, \( \nabla_\mu \) denotes the covariant derivative associated with the gravitational metric \( g_{\mu\nu} \).

We assume that the Lagrangian \( L(X, \phi) \) is a general function satisfying the following conditions: \( L_X \geq 0 \) (Null Energy Condition) and \( 2X L_{XX}/L_X > -1 \) (Hyperbolicity condition). The first condition guarantees that the perturbations carry positive kinetic energy while the second one implies the stability with respect to high frequency perturbations, and is necessary for the Cauchy problem to be well posed (see e.g. \([60, 62, 63, 31, 20, 102]\)). These conditions restrict the variety of the allowed Lagrangians along with the corresponding solutions, and are unavoidable for any physically meaningful model. The energy-momentum tensor of the theory is

\[ T_{\mu\nu} = L_X \nabla_\mu \phi \nabla_\nu \phi - g_{\mu\nu} L. \]  

It is well known (see e.g. \([32]\)), that for time-like derivatives (i.e., \( X > 0 \)), the models under consideration can be described in a hydrodynamical language by introducing an effective four velocity

\[ u_\mu = \frac{\nabla_\mu \phi}{\sqrt{2X}}, \]  

along with the pressure

\[ p = L(X, \phi), \]

\[ ^4 \text{We use the notation } (\ldots)_X \equiv \frac{\partial (\ldots)}{\partial X} \text{ and the signature (+ −−−).} \]

\[ ^5 \text{For further discussion of these topics please refer to appendix §A.} \]

\[ ^6 \text{For a different opinion see [103].} \]

\[ ^7 \text{Note that even for } X > 0 \text{ the effective four velocity introduced in (4.2) is not necessarily future directed. However, the analogy with the perfect fluid can be made exact by multiplying this expression (4.2) with } \pm 1 \text{ so that } u^0 > 0. \]
the energy density

\[(4.4)\]
\[
\rho(X, \phi) = 2XL_X - L,
\]

and the sound speed\(^8\)

\[(4.5)\]
\[
\epsilon_s^2(X, \phi) = \left(1 + \frac{2XL_X}{L_X}\right)^{-1} = \left(\frac{\partial p}{\partial \rho}\right)_\phi.
\]

In these variables the energy-momentum tensor has the form corresponding to the one of a perfect fluid:

\[
T_{\mu\nu} = (\rho + p)u_\mu u_\nu - pg_{\mu\nu}.
\]

It is convenient to use the hydrodynamical notation for these functions of \(\phi\) and \(X\) also for \(X \leq 0\) when they do not have their usual physical meaning of velocity etc.

### 4.2.1 Field redefinitions and conditions for stationarity

If the field \(\phi\) does not have any direct interactions except with gravity, then obviously a field redefinition \(\phi = \phi(\tilde{\phi})\) cannot affect any observables besides the field itself. This is a particular case of a stronger statement (see e.g., [104]). Obviously the solutions \(\phi(x)\) and \(\tilde{\phi}(x)\) result through Einstein equations in the same gravitational metric \(g_{\mu\nu}(x)\), and describe in that sense the same physical process. Thus, it is interesting to investigate the properties of k-essence under field redefinitions. Under the field redefinition \(\phi = \phi(\tilde{\phi})\), we have \(\nabla_\mu \phi = \left(d\phi/d\tilde{\phi}\right)\nabla_\mu \tilde{\phi}\); whereas the expressions for the energy-momentum tensor \(T_{\mu\nu}\), and all hydrodynamical quantities \(\rho, p, c_s\) and \(u^\mu\) remain unchanged or covariant\(^9\). Here we should distinguish between covariance and invariance. Covariance means that the way how the quantities/equations are constructed from other objects remains unchanged whereas invariance implies exactly the same functional dependence on these objects. For example, the formula

\(^8\)This formula for the sound speed was introduced for the cosmological perturbations in [34]. One can show [31] that the same expression is valid in the general case of backgrounds with time-like field derivatives: \(X > 0\).

\(^9\)Note that the four velocity (4.2) is invariant up to the sign only.
defining the energy density \( \rho \) through Lagrangian \( L, X \) and the derivative \( L_X \) looks the same after a field redefinition (covariant); however, the dependence of the Lagrangian on the field does change (not invariant). It is obvious that, for example, the value of physical energy density at every point should not change under field redefinitions, but here these quantities reveal in addition such covariance with respect to the field redefinitions as it is the case, for example, the Euler-Lagrange equations. However, this covariance is not guarantied for all interesting objects. It is worthwhile mentioning that, for example, the metric \[ G_{\mu\nu} [\phi_0] = \frac{L_X}{c_s} \left( g_{\mu\nu} - \frac{c_s^2 L_{XX}}{L_X} \nabla_{\mu} \phi_0 \nabla_{\nu} \phi_0 \right), \]
describing the propagation of small perturbations \( \pi \) around a given background \( \phi_0 (x) \), transforms conformally under field redefinitions \( \phi = \phi (\tilde{\phi}) \):

\[
G_{\mu\nu} [\phi_0] = \left( \frac{d\tilde{\phi}}{d\phi} \right)_0^2 G_{\mu\nu} [\tilde{\phi}_0].
\]

Thus, as expected, the causal structure does not change under field redefinitions. The conformal factor \( \left( d\tilde{\phi}/d\phi \right)_0^2 \) compensates for the redefinition of perturbations \( \pi = \left( d\phi/d\tilde{\phi} \right)_0 \tilde{\pi} \).

Let us further consider a stationary space-time with metric \( g_{\mu\nu} \) and a time-like Killing vector \( t^\alpha \). Thus, \( \mathcal{L}_t g_{\mu\nu} = 0 \), where \( \mathcal{L}_t \) is the Lie derivative. The configuration is stationary if by definition

\[
\mathcal{L}_t T_{\mu\nu} = 0.
\]

Using Leibniz’s rule we have

\[
\mathcal{L}_t T_{\mu\nu} = (\mathcal{L}_t L_X) \nabla_\mu \phi \nabla_\nu \phi - g_{\mu\nu} \mathcal{L}_t L
\]

\[
+ L_X \left[ (\mathcal{L}_t \nabla_\mu \phi) \nabla_\nu \phi + (\mathcal{L}_t \nabla_\nu \phi) \nabla_\mu \phi \right] = 0.
\]

(4.6)
By multiplying this expression with \( g^{\mu\nu} \) we obtain

\[
(4.7) \quad 0 = \mathcal{L}_t T^\mu_\mu = \mathcal{L}_t (2XL_X - 4L) = \mathcal{L}_t (\rho - 3p).
\]

Suppose the configuration \( \phi(x^\mu) \) is such that \( \nabla_\mu \phi \) is a null vector: i.e., \( X = 0 \). In that case we can multiply the right hand side of the equation (4.6) with \( g^{\mu\nu} \) to obtain \( \mathcal{L}_t L = 0 \). Further, we have \( 0 = \mathcal{L}_t L = L_\phi \partial_t \phi \). As we are looking for stationary but not static solutions, we have \( L_\phi = 0 \). Thus, the Lagrangian should be symmetric with respect to field shifts \( \phi \to \phi + c \), where \( c \) is an arbitrary constant.

For \( X \neq 0 \) it is convenient to introduce the projector

\[
(4.8) \quad \mathcal{P}_{\mu\nu} = g_{\mu\nu} - \frac{\nabla_\mu \phi \nabla_\nu \phi}{2X},
\]

with the properties

\[
(4.9) \quad \mathcal{P}_{\mu\nu} \nabla^\nu \phi = 0, \quad \mathcal{P}_{\mu\lambda} \mathcal{P}^{\lambda\nu} = \mathcal{P}_{\mu}^{\nu}, \quad \text{and} \quad \mathcal{P}_{\mu}^{\mu} = 3.
\]

Moreover, this projector is both invariant and covariant under field reparametrizations: \( \mathcal{P}_{\mu\nu}[\phi] = \mathcal{P}_{\mu\nu}[\tilde{\phi}] \). By acting with the projector \( \mathcal{P}_{\mu\nu} \) on the left hand side of equation (4.6), we have \( 0 = \mathcal{P}_{\mu\nu} \mathcal{L}_t T^\mu_\nu = -3\mathcal{L}_t L \). Therefore, if the configuration is stationary then in particular

\[
(4.10) \quad \mathcal{L}_t L = 0,
\]

which for the hydrodynamical case reduces to the constancy of the pressure \( p \). Combining this with (4.7) we obtain the time independence of the energy density \( \rho \) or

\[
(4.11) \quad \mathcal{L}_t (XL_X) = 0.
\]

Further, we can act on the left hand side of equation (4.6) with \( \mathcal{P}^{\alpha\nu} \) so that

\[
0 = \mathcal{P}^{\alpha\nu} \mathcal{L}_t T^\mu_\nu = L_X \mathcal{P}^{\alpha\nu} (\mathcal{L}_t \nabla_\nu \phi) \nabla_\mu \phi.
\]

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Thus, stationarity implies

\[ \mathcal{P}^{\alpha\beta} (\mathcal{L}_t \nabla_\beta \phi) = 0. \]

Using the properties of the projector (4.9), Leibniz’s rule and that \( t^\alpha \) is a Killing vector, one obtains

\[ 0 = \mathcal{P}^{\alpha\beta} (\mathcal{L}_t \nabla_\beta \phi) = -\nabla_\beta \phi \mathcal{L}_t \mathcal{P}^{\alpha\beta} = -\mathcal{L}_t \nabla^\alpha \phi + \frac{\nabla^\alpha \phi}{2X} \mathcal{L}_t X. \]

The last expression in turn can be written in the following form

\[ -\mathcal{L}_t \nabla^\alpha \phi + \frac{\nabla^\alpha \phi}{2X} \mathcal{L}_t X = \sqrt{2X} \mathcal{L}_t \left( \frac{\nabla^\alpha \phi}{\sqrt{2X}} \right). \]

Therefore, stationarity implies

(4.12) \[ \mathcal{L}_t \left( \frac{\nabla^\alpha \phi}{\sqrt{2X}} \right) = 0, \]

or in the hydrodynamical notation \( \mathcal{L}_t u^\mu = 0 \)^1. Thus, we have proved that for any stationary configuration the following conditions

(4.13) \[ \mathcal{L}_t u^\mu = 0, \quad \mathcal{L}_t \rho = 0, \quad \text{and} \quad \mathcal{L}_t p = 0, \]

should be satisfied. Note that these conditions are covariant under field redefinitions and, for the hydrodynamical case (\( X > 0 \)), are intuitively clear requirements. Sometimes (see e.g., [94]) one claims that the stationarity implies a stronger requirement:

(4.14) \[ \mathcal{L}_t \nabla_\mu \phi = 0, \]

instead of the condition (4.12). However, the equation above is not covariant under field redefinitions, and does not follow from the stationarity of the energy-momentum tensor.

^1The vector \( u^\mu \) is formally imaginary for \( X < 0 \). However, without any change of the results one could redefine \( u^\mu \) in this case: \( u^\mu = \nabla^\mu \phi / \sqrt{-2X} \).
Now let us find what type of theories \( L(X, \phi) \) and field configurations \( \phi(x^\mu) \) can, in principle, satisfy conditions (4.13). It is convenient to choose a coordinate system \((t, x^i)\) such that the time coordinate corresponds to the integral curves of \( t^\alpha \). In that case the Lie derivative reduces to the partial derivative \( \mathcal{L}_t = \partial_t \).

4.2.2 Which field configurations can have constant effective four velocity \( u^\mu \)?

Now let us find the configurations \( \phi(x^\mu) \) satisfying the condition on the effective four velocity (4.12). For the time component of the four velocity we have

\[
\partial_t \left( \frac{\dot{\phi}}{\sqrt{2X}} \right) = \frac{\ddot{\phi}}{\sqrt{2X}} + \dddot{\phi} \partial_t \left( \frac{1}{\sqrt{2X}} \right) = 0,
\]

where \( \dot{\phi} = \partial_t \phi \), while for the spatial components

\[
\partial_t \left( \frac{\partial_i \phi}{\sqrt{2X}} \right) = \frac{\partial_i \dot{\phi}}{\sqrt{2X}} + \partial_i \phi \partial_t \left( \frac{1}{\sqrt{2X}} \right) = 0.
\]

Obviously these equations have a trivial static solution \( \phi = \phi(x^i) \). To find a nontrivial solution we combine these two equations to obtain following system of equations:

\[
\dot{\phi} \partial_i \dot{\phi} - \dddot{\phi} \partial_t \phi = 0,
\]

which is equivalent to

\[
\partial_t \left( \frac{\partial_i \phi}{\dot{\phi}} \right) = 0.
\]

This is a system of second order partial differential equations. Integrating equation (4.17) we obtain the following linear homogeneous system

\[
\partial_i \phi = V_i (x^j) \dot{\phi},
\]

where \( V_i (x^j) \) are unknown time independent functions. This is a first order system of three partial differential equations for only one function \( \phi \). Let us find the consistency
conditions under which the system can have solutions. Differentiating the $i-$equation with respect to $x^j$, and using the time differentiation of the $j-$equation we obtain
\[ \partial_j \partial_i \phi = \partial_j V_i \dot{\phi} + V_i \partial_j \dot{\phi} = \partial_j V_i \dot{\phi} + V_i V_j \ddot{\phi}. \]

Now we can compare this result with the result of the same procedure performed for the $j-$equation. We obtain
\[ \partial_i V_j - \partial_j V_i = 0. \]

For a simply connected manifold, the last equation implies the existence of a function (potential) $\Psi \left( x^i \right)$ such that $V_i = \partial_i \Psi$. Otherwise, there are no solutions for (4.18).

For the $i-$equation we can assume that all $x^k$ with $k \neq i$ are frozen parameters, and for the characteristics (for the method of characteristics see, for example, the excellent book [105]) we obtain
\[ \frac{dt}{d\tau} = -\partial_i \Psi \left( x^j \right), \quad \frac{dx^i}{d\tau} = 1. \]

The first integral $I$ of this system is given by the constant of integration for the equation
\[ \frac{dt}{dx^i} = -\partial_i \Psi \left( x^j \right). \]

By integrating this we obtain
\[ t = I - \Psi \left( x^i \right). \]

Therefore, the general solution $\phi \left( t, x^i \right)$ is given as an arbitrary function of the first integral $I$:
\[ (4.19) \quad \phi \left( t, x^i \right) = \Phi \left( t + \Psi \left( x^i \right) \right). \]
Thus, the general solution for equations (4.15) and (4.16) contains two arbitrary functions. Note that the system (4.17) does not have any other general solutions besides (4.19). It is easy to prove that this solution satisfies the equations (4.15) and (4.16). Indeed, we have

\[ \dot{\phi} = \frac{d\Phi}{dI}, \quad \text{and} \quad \partial_i \phi = \frac{d\Phi}{dI} \partial_i \Psi. \]

Therefore,

\[ (4.20) \quad X = \frac{1}{2} \left( \frac{d\Phi}{dI} \right)^2 \left( g^{00} + 2g^{0i} \partial_i \Psi + g^{ik} \partial_i \partial_k \Psi \right), \]

and the time component

\[ \frac{\dot{\phi}}{\sqrt{2X}} = \frac{1}{\sqrt{g^{00} + 2g^{0i} \partial_i \Psi + g^{ik} \partial_i \partial_k \Psi}}, \]

along with the spatial components

\[ \frac{\partial_i \phi}{\sqrt{2X}} = \frac{\partial_i \Psi}{\sqrt{g^{00} + 2g^{0i} \partial_i \Psi + g^{ik} \partial_i \partial_k \Psi}}, \]

are obviously time independent because the metric is stationary. It is worth mentioning that by using the condition (4.14) we would arrive at the general solution \( \phi(t, x^i) = t + \Psi(x^i) \); missing the arbitrary functional dependence \( \Phi \). Note that arbitrary field redefinitions correspond to the freedom in choosing \( \Phi \).

4.2.3 Which Lagrangians Allow for Stationary Configurations?

Now let us consider the restrictions on \( L(X, \phi) \) arising from the requirement that the pressure and energy density should be time independent for the general solution (4.19). From equation (4.10), we have

\[ (4.21) \quad \partial_t L = L_{\phi} \dot{\phi} + L_X \dot{X} = 0, \]

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while from equation (4.11)

\[ \partial_t (XL_X) = \dot{XL}_X + XL_{X\phi} \dot{\phi} + XL_{XX} \ddot{X} = 0. \]

Eliminating \( \dot{X} \) from these equations results in

\[ XL_{X\phi} - (XL_{XX} + L_X) \frac{L_{\phi}}{L_X} = 0. \] (4.22)

This equation is a second order partial differential equation for \( L(X, \phi) \). A trivial solution of this equation is a shift symmetric Lagrangian \( L(X) \). It is well known that shift symmetric theories are exactly equivalent to hydrodynamics for \( X > 0 \). Obviously hydrodynamics allows for steady flows. Let us find a general solution of the equation (4.22). This general solution should depend on two arbitrary functions.

It is convenient to rewrite equation (4.22) in the following form:

\[ \frac{\partial \ln (L_{\phi}/L_X)}{\partial \ln X} = 1. \]

Integrating this equation we obtain

\[ L_{\phi} = \sigma (\phi) XL_X, \] (4.23)

where \( \sigma (\phi) \) is an arbitrary function. The last equation (4.24) is a first order linear partial differential equation. Similar to our previous calculations, we use the method of characteristics to find the general solution. For the characteristics we have

\[ \frac{d\phi}{d\tau} = 1, \quad \text{and} \quad \frac{dX}{d\tau} = -\sigma (\phi) X. \] (4.24)

Thus, the integral curves are given by the equation

\[ \frac{dX}{d\phi} = -\sigma (\phi) X. \]

The general solution of the last equation is

\[ X = \mathcal{J} \exp \left( - \int \sigma (\phi) \, d\phi \right), \]
where $\mathcal{J}$ is a constant of integration. Thus, the general solution to the equations (4.23) and (4.22) is an arbitrary function of the first integral $\mathcal{J}$ of the dynamical system (4.24):

\begin{equation}
L(X, \phi) = F(X e^{f(\phi)}),
\end{equation}

where $F$ and $f(\phi) = \int \sigma(\phi) \, d\phi$ are arbitrary functions. Note that all solutions of (4.22) are described by (4.25). It is obvious that the Lagrangian (4.25) has a hidden shift symmetry. Namely, we can always perform the field redefinition

\begin{equation}
\bar{\phi}(\phi) = \int d\phi \, e^{f(\phi)/2},
\end{equation}

so that the new Lagrangian is shift symmetric: $L(X, \phi) = F(\bar{X})$, where $\bar{X} = \frac{1}{2} g^{\mu\nu} \nabla_\mu \bar{\phi} \nabla_\nu \bar{\phi} = X e^{f(\phi)}$. Thus, all scalar field theories that allow for stationary configurations are necessarily shift symmetric (explicitly or after a field redefinition).

Further, we will use the notation $\bar{\phi}$ always for such field variables in which the system is invariant under shift transformations: $\bar{\phi} \to \bar{\phi} + c$, where $c$ is an arbitrary constant.

Finally, we can specify the profiles $\Phi$ of stationary configurations. Equations (4.21) and (4.20) yield

\begin{equation}
L_{\phi} + L_X \left( \frac{d^2 \Phi}{dT^2} \right) \left( g^{00} + 2 g^{0i} \partial_i \Psi + g^{jk} \partial_j \partial_k \Psi \right) = 0,
\end{equation}

and using equations (4.23) and (4.20) we obtain

\begin{equation}
\frac{d^2 \Phi}{dT^2} + \frac{1}{2} \left( \frac{d\Phi}{dT} \right)^2 \sigma(\Phi) = 0.
\end{equation}

We know that in terms of the new field $\bar{\phi}$, the Lagrangian is shift symmetric. Thus, for this parametrization $\sigma(\bar{\phi}) = 0$. Therefore, $\Phi(I) = \alpha I + \beta = t + \Psi(x^i)$ where we have absorbed the constants into $\Psi$ and $t$. Thus, in terms of the field variable $\bar{\phi}$, in which the theory is shift symmetric, the possible stationary configurations are
always given by

\begin{equation}
\tilde{\phi} = t + \Psi (x^i),
\end{equation}

and we are back to the usual ansatz (4.14). The stationary configurations in terms of
the field variable $\phi$ can be obtained by solving equation (4.26) or (4.27) with respect
to $\phi$. This procedure determines the function $\Phi$, while the function $\Psi (x^i)$ has to be
fixed from the equations of motion and boundary/initial conditions.

It is worth noting that if the metric $g_{\mu \nu}$ possesses another Killing vector corre-
sponding to, e.g., axial symmetry: $\mathcal{L}_\theta g_{\mu \nu} = 0$, then we can apply the result (4.28) to
the angular variable $\theta$. Thus, the solution is

\begin{equation}
\tilde{\phi} = t + \Omega \theta + \Psi (x^i_\perp),
\end{equation}

where $\Omega$ is a constant and $x^i_\perp$ denotes the rest of the coordinates.

4.3 Conclusion

In this chapter we have proved that the existence of stationary configurations
requires shift symmetry. Namely (may be after a field redefinition) the system has
to be invariant with respect to the transformation $\tilde{\phi} \rightarrow \tilde{\phi} + c$, for all constants $c$.

The result is valid in the self-consistent case where the geometry is produced by
the scalar field, as well as in the test field approximation where the stationary field
configuration appears on the gravitational background governed by other sources.

The shift symmetry implies the conservation of the Noether current

\[ J_\mu = L_{\tilde{X}} \nabla_\mu \tilde{\phi}. \]

Interestingly, the equation of motion implies $\nabla_\mu J^\mu = 0$, which is a statement of the
conservation of the current $J^\mu$. In the case when $\nabla_\mu \tilde{\phi}$ is time-like, the current $J^\mu$
can be written in the form of an effective particle density current \( J^\mu = \tilde{n} u^\mu \), where the particle density\(^{11}\) is

\[
\tilde{n} = \sqrt{2\, \tilde{\mathcal{X}} L_{\tilde{\mathcal{X}}}}.
\]

Note that this current is not covariant under field redefinitions. The conservation of the particle density current usually holds in standard hydrodynamics. However, the most interesting models of cosmological scalar fields do not possess this additional conservation law associated with shift symmetry. Thus, the result obtained in this chapter implies that there is no exact Bondi (steady flow) accretion for popular classes of models of dynamical dark energy like quintessence and k-essence. This result may not have a very strong qualitative impact on the growth of black holes or on the evolution of the cosmological fields around them. Indeed, one should expect that the accretion rate should be in any case rather small (for the case of canonical scalars see [95]). Especially in the late/current universe, one can almost always neglect the growth of the black hole along with the corresponding back-reaction. Nevertheless, this result changes the setup for the investigation of the problem. Now in order to study how these fields could accrete onto black holes, one is forced to solve the Cauchy problem for nonlinear partial differential equations, instead of solving the boundary problem for nonlinear ordinary differential equations. In particular, to approach this problem one has to choose some initial configuration for the field and its time derivative. At this stage, it is not clear what are reasonable, physically motivated initial conditions, and at what time they should be posed. This is very different from the case of Bondi accretion where the boundary conditions are fixed by cosmological evolution, and the membrane property of the BH horizon. However, it may happen that there are some special attractor or self-similar regimes to which

\(^{11}\)Note that this number density is none other than the canonical momenta for the field \( \tilde{\phi} \) in the comoving reference frame.
the solutions would approach at late time. Nevertheless, one cannot guarantee either
the existence of these attractors, nor their uniqueness for a general model. Moreover,
even if a unique attractor exists, then it is not a priori known how wide the basin
of attraction is in the phase space consisting of initial configurations of the field and
its time derivative. Thus, the procedure for finding these attractor solutions is not
only a predominantly numerical exercise, but also generically not very promising and
predictive. Nevertheless, it is very interesting to find examples of scalar field systems
possessing solutions of this type. In [95] it was demonstrated that for canonical
scalars and many potentials the solutions indeed approach a steady flow.

In addition, one has to mention that having a shift symmetric theory is a necessary,
but insufficient condition for the existence of stationary configurations. For exam-
ple, in hydrodynamics there can be either exceptional theories or even exceptional
boundary conditions for which there are no stationary configurations. In particular
the simple accretion of dust onto a black hole occurs along geodesics and therefore is
not steady. A similar situation happens in the case of the ghost condensate for which
the accretion rate blows up when the field configuration at spatial infinity approaches
the condensation point (compare [94] and [92]). Moreover, in the DBI model consid-
ered in [33], it was found that a physically meaningful steady state accretion is not
possible when the sound speed at spatial infinity is $c_s^2 > 4/3$.

In this chapter we have considered only a single self interacting scalar field. It
would be interesting to study other types of fields, in particular one could think of
scalars with internal degrees of freedom, e.g., charged scalars accreting onto a charged
black hole. We expect that the appearance of new external forces and internal degrees
of freedom can change the picture. Another interesting problem is to find possible
attractor or self-similar asymptotic solutions, and develop a perturbation theory
around them. As we have shown, stationary configurations are possible only for theories that are equivalent to perfect fluids. This result reveals once again that the relation between hydrodynamics and field theory is rather deep. Therefore, we think this connection deserves a further study. We found that investigation of possible dynamical backgrounds around black holes is interesting not only from the point of view of mathematical physics, but may be relevant for a better understanding of both black holes physics and may be even the nature of dark energy.
CHAPTER V

Reconstruction of Non-Canonical Inflationary Actions

5.1 Introduction

Since the landmark COBE experiment, the study of the cosmos has entered a new age of precision cosmology. For the first time in the history of modern cosmology, direct quantitative measurements of early cosmological observables are available. The data taken by COBE was critical in establishing inflation as the central paradigm in our theories of the origin of the universe [7]. Thanks to experiments that measured the spectrum of CMBR fluctuations, we have now confirmed that the near-scale invariance of large scale fluctuations that is a prediction of inflation are in fact borne out in the data. Although observation supports the general theory of inflation, as of now the data is unable to conclusively determine the mechanism responsible for inflation.

The difficulty in discriminating between different inflation models lies in the fact that all current models of inflation predict the same near-scale invariant spectrum of scalar fluctuations. To further narrow down the number of observationally consistent inflation models, observables independent of the scalar perturbation need to be measured and compared to model predictions. Two additional inflationary observables are the spectrum of tensor perturbations $P_h$ [106, 107] and the non-
gaussianity $f_{NL}$ [35, 37, 38] of the CMBR temperature perturbation spectrum. In recent years greater progress has been made in measuring these quantities directly. Upper bounds on the amplitude of the tensor perturbation spectrum, which is the spectrum of relic gravitational waves, have been determined directly through analysis of the CMBR polarization [108, 109]. The non-gaussianity, which represents the deviation of the curvature perturbation from gaussian statistics, is also being better understood. Analysis of WMAP3 data [110] has found evidence of non-gaussian statistics in the CMBR temperature spectrum. With a better knowledge of these extra observables it becomes possible to better determine which model of inflation is most likely to have taken place. For example, a large $f_{NL}$ would tend to rule out a single field inflation model with minimal kinetic terms, while favoring those models that predict a large non-gaussianity.

Ultimately, one would like to use the features of the CMBR temperature anisotropy to reconstruct the inflaton action directly. It is customary to write the general scalar field action as [34]

$$S = \int d^4x \sqrt{-g} L(X, \phi), \quad X = \frac{1}{2} g_{\mu\nu} \partial_\mu \phi \partial_\nu \phi.$$  

In our analysis we will limit ourselves to actions that contain no third or higher derivatives of the inflaton. Throughout this chapter we will assume that the curvature of the three non-compact space dimensions will be zero. Following the cosmological principle, we use the FRW metric: $g_{\mu\nu} = \text{diag}(1, -a^2, -a^2, -a^2)$, where $a$ is a time dependent scale factor. Using (5.1), the Friedmann equations for the scale factor are

$$3M_{pl}^2 H^2 = \rho,$$

$$-2M_{pl}^2 \dot{H} = \rho + p,$$

where $H = \frac{d \log a}{dt}$ is the Hubble parameter. The quantities $p$ and $\rho$ are the pressure
and energy density, respectively, of the scalar field, which in terms of the Lagrangian are

\begin{equation}
 p = L, \quad \rho = 2XL_X - L.
\end{equation}

(5.4)

In single field inflation models with a minimal kinetic term the action is

\begin{equation}
 S = \int d^4 x \sqrt{-g} [X - V(\phi)].
\end{equation}

(5.5)

If we assume (5.5), the only function that needs to be determined from the data is the potential \( V(\phi) \). Reconstruction of the inflationary potential for models of the form (5.5) has been studied extensively [111]. However, by assuming that the action has a minimal kinetic term we neglect a rich class of models such as DBI inflation [28, 57], k-inflation [32] and ghost inflation [112]. In contrast, only a handful of articles have been written that deal with the reconstruction of inflationary actions with general kinetic terms [36, 113, 114].

In non-minimal kinetic models the speed at which scalar fluctuations propagate can be different than the speed of light. This can affect the temperature anisotropy in two ways. First, if \( c_s < c = 1 \), scalar fluctuations have a sound horizon that is smaller than the cosmological horizon, causing curvature perturbations to freeze in earlier than normal. Depending on how the Hubble parameter and sound speed change during the course of inflation, the temperature anisotropy can develop noticeable signatures of non-minimal kinetic terms. Second, models with non-minimal kinetic terms will in general produce a non-gaussian spectrum. Traditionally, the non-gaussianity is measured by the nonlinearity parameter \( f_{NL} \) defined by the following ansatz for the curvature perturbation:

\begin{equation}
 \zeta = \zeta_L - \frac{3}{5} f_{NL} \zeta_L^2.
\end{equation}

(5.6)
Here, $\zeta$ is the general curvature perturbation and $\zeta_L$ is a curvature perturbation with gaussian statistics. Within the standard canonical action (5.5), non-gaussianities can be produced by cubic or higher order terms in the inflaton potential, or by secondary interactions with gravity [38]. However, non-gaussianities produced by these mechanisms are on the order of the slow roll parameters, and thus small. In contrast, models with non-minimal kinetic terms can have large non-gaussianities, providing a clear distinction from canonical inflation.

The goal of this chapter will be to reconstruct an inflationary action from observables starting with as few initial assumptions as possible. In this chapter we take the experimental inputs to be the scalar curvature perturbation $P_s$, the tensor curvature perturbation $P_h$, and the non-gaussianity (nonlinearity) parameter $f_{NL}$.

Unfortunately, completely reconstructing the off-shell action is not possible since the observables only carry on-shell information. To understand why the off-shell action is inaccessible to us, consider the interpretation of the Lagrangian $L(X, \phi)$ as a surface in the three dimensional space $(\phi, X, L)$ [114]. Because the observables are insensitive to the off-shell behavior of the Lagrangian, we can only determine the one-dimensional trajectory $L = L(X(\phi), \phi)$ of the on-shell Lagrangian, embedded in the two-dimensional surface defined by $L = L(X, \phi)$. A one-dimensional trajectory has an infinite number of surfaces that contain it, each related to one another by a canonical transformation [114]. Therefore we have to be more specific about the form of the Lagrangian that we are trying to find. In this chapter we will reconstruct inflationary Lagrangians that have the form

$$L(X, \phi) = P(g_1(X), \ldots, g_m(X), f_1(\phi), \ldots, f_n(\phi)). \tag{5.7}$$

Here it is assumed that $P(x_1, \ldots, x_m, y_1, \ldots, y_n)$, which we will refer to as the partition of the action, is a known function of the $\{x_i\}$ and $\{y_\alpha\}$, and the functions $\{g_i\}$ and
\{f_\alpha\} are not all known. Once the on-shell trajectory \( \phi = \phi(k) \) is determined, the Lagrangian \((5.7)\) defines a surface in the \((\phi, X, L)\)-space up to a field redefinition.

Before non-canonical inflation models were first considered, reconstructions of the inflationary action always assumed that the Lagrangian had the form

\[ L(X, \phi) = X - V(\phi). \]

Here the dependence of the action on \( X \) is known, and the problem of reconstructing the inflaton action is reduced to finding the inflaton potential \( V(\phi) \) from the data. In our language, the partition of the action that was assumed was

\[ P(g(X), f(\phi)) = g(X) - f(\phi). \]  \hspace{1cm} (5.8)

Since the scalar field action is assumed to be canonical, then the function \( g(X) \) is taken as a known and is simply \( g(X) = X \). The function \( f(\phi) \) is the unknown and represents the potential of the canonical scalar field that previous inflationary action reconstructions were concerned with. The procedure that we develop here is a generalization of procedures used to determine the potential in canonical scalar field models of inflation. For example, in the case of the partition \((5.8)\), our procedure does not require that any assumptions be made about the form of \( g(X) \). Instead, \( g(X) \) and \( f(\phi) \) are treated on an equal footing, and our procedure can determine both using CMBR data.

The idea will be to use data on the CMBR perturbation spectrum to find the functions \( \{g_i(X)\} \) and \( \{f_\alpha(\phi)\} \). In a naive comparison with algebraic linear equations, we expect that if there are \( n \) unknown functions, finding the action requires \( n \) experimental inputs. Since we are assuming that there are only three observables: \( P_s, P_h \) and \( f_{NL} \), we can derive three reconstruction equations, which can determine an action with three or fewer unknown functions. In the case where the number of
unknown functions is less than the number of experimental inputs, the reconstruction
equations not used to find the action become constraint equations.

Since we are interested in solving for functions \( \{g_i\} \) and \( \{f_\alpha\} \) and not just numbers,
we will need to know at least a portion of \( P_s, P_h \) and \( f_{NL} \) as functions of the scale \( k \). While the scale dependence of \( P_s \) is known to be at least approximately power-law dependent on \( k \), the scale dependence of the other two observables \( P_h \) and \( f_{NL} \) is unclear at this point. Although future experiments will be able to clarify some aspects of the tensor and non-gaussianity signals, their exact functional forms will probably not be available for quite some time if at all. Regardless, the method we develop here does have utility outside of reconstruction. This method is well suited to testing how the form of an action depends on the observables. For instance if the scalar perturbation is of the near-scale invariant variety:

\[
P_s \propto k^{n_s-1},
\]

(5.9)

we can use \( P_s \) to help derive an action and study its dependence on the index \( n_s \). That way if we wish to connect the action derived from (5.9) to an action derived from theory, we can see if the theoretical action leads to reasonable results for the observables. Furthermore, as we mentioned earlier when there are only one or two unknown functions in (5.7), the remaining reconstruction equations determine new consistency relations. In this chapter most of the examples we deal with have only two unknown functions, which we solve for using the scalar and tensor spectrum data. The reconstruction equation derived from the non-gaussianity will then be a constraint; relating the non-gaussianity to the sound speed, the Hubble parameter and/or their derivatives. Outside of deriving the action, we also find a method for quickly obtaining the sound speed as a function of time from the scalar and tensor spectrum

\footnote{Here, \( f_{NL} \) represents the equilateral bispectrum, and is therefore a function of a single scale.}
perturbation spectra. Finding a sound speed different from the speed of light even over a small range of scales would be a powerful indication of non-canonical inflation.

This chapter is organized as follows. In section §5.2 we present the method for reconstructing the action from the scalar, tensor and non-gaussianity spectra. We explain how cosmological data can find the Hubble parameter $H$, and the sound speed $c_s$, and how these in turn can be used to find three unknown functions of the action (5.7). Once the method has been explained in section §5.2.1 we carry out a derivation of the action for different functions $P(z_1, z_2, z_3)$, assuming that the scalar and tensor power spectra both scale like $k$ to some power. In section §5.3 we apply our method to find the warp factor and potential in a generalized DBI inflation model. We find the warp factor and potential as functions of the spectral indices and the initial value of the Hubble parameter. The results for these are compared to the theoretically motivated warp factor and potential used in D3 brane DBI inflation. Finally, in section §5.4 we review our main results.

5.2 The Reconstruction Equations

We start our derivation of the reconstruction procedure by explaining how the observables are used to find the Hubble parameter $H$ and sound speed $c_s^2$. Once we have these, the action can be obtained using a set of reconstruction equations that will be shown later. Let us begin by recalling the definition of the slow roll parameter$^3$ $\epsilon$ in terms of the Hubble parameter. The definition of $\epsilon$ implies that

\[
\frac{dH}{dt} = -\epsilon H^2.
\]  
(5.10)

Since the perturbation spectra and non-gaussianity are functions of $k$ and not time, we wish to rewrite this equation for $\frac{dH}{dt}$ into an equation for $\ddot{H} = \frac{dH}{d\log k}$. However,

$^3$The term “slow roll parameter” is taken from chaotic inflation where inflation occurs only when the inflaton “velocity” $\dot{\phi}$ is small. However, DBI inflation can still occur for large $\phi$. 

$^2$The method described here was inspired by the technique used in [115]
because we are assuming a general sound speed, we must be careful to differentiate between the horizons of scalar and tensor fluctuations. If the sound speed differs from unity (in particular $c_s < 1$), then the horizon size of scalar fluctuations: $(aH/c_s)^{-1}$ is smaller than that of the tensor fluctuations: $(aH)^{-1}$. This implies that at any given time, the scales $k_s$ and $k_t$ at which the scalar and tensor fluctuations leave their respective horizons, will in general be different. For our purposes, we choose to study the dependence of $H$ on the scalar wave number $k_s$. Therefore, the condition for horizon exit is now $kc_s = aH$ instead of the more familiar relation: $k = aH$.

Having made clear our choice of wave number, we now set out to express $\frac{d \log k}{dt}$ in terms of familiar quantities:

$$
\frac{d \log k}{dt} = H(1 - \epsilon + \kappa \frac{d \log k}{dt}),
$$

where we have defined $\kappa = -\frac{c_s}{c_s}$. Solving for $\frac{d \log k}{dt}$ above:

$$(5.11) \quad \frac{d \log k}{dt} = \frac{H(1 - \epsilon)}{1 - \kappa}.
$$

With equation (5.11) in hand, the equation for $\ddot{H}$ can now be found from (5.10):

$$(5.12) \quad \ddot{H} = -\frac{H\epsilon(1 - \kappa)}{1 - \epsilon}.
$$

We will use this equation to find $c_s$ once we have found $H$ and $\epsilon$ in terms of the observables. Since $\epsilon$ depends on the time derivative of the Hubble parameter, $H(k)$ and $\epsilon(k)$ are independent parameters. Since we have two independent parameters, we will likely need two independent observables. The two observables we will use here are the scalar and tensor perturbation spectra. Recall that to first order in the slow roll parameters the perturbation spectra are given by

$$(5.13) \quad P_s(k) = \frac{H^2}{8\pi^2 M_{pl}^2 c_s} \bigg|_{k_s c_s = aH},$$
The extra source of information can also be garnered from the non-gaussianity parameter $f_{NL}$. However, going in this route would result in a more complicated solution. The parameter $\epsilon$ can be found as a function of wave number using (5.13). Solving for $\epsilon$ we obtain

$$
\epsilon|_{k_s c_s = aH} = \frac{1}{8\pi^2 M_{pl}^2 P_s(k_s)} \left. \frac{H^2}{c_s} \right|_{k_s c_s = aH}.
$$

(5.15)

As a matter of convenience define $P_s = A P_s$, where $A$ is the value of the scalar perturbation at some fiducial scale $k = k_0$. If $k_0 \simeq 0.002 \text{Mpc}^{-1}$ then present data suggests that $A \simeq 10^{-9}$. Here, $P_s$ is the normalized scalar perturbation defined such that $P_s(k_0) = 1$. Furthermore, let $H = \mathcal{A} \mathcal{H}$ where $A^2 = 8\pi^2 M_{pl}^2 A$. Substituting (5.15) in for $\epsilon$ in equation (5.12) we have

$$
\left. \mathcal{H} \right|_{k_s c_s = aH} = -\frac{\mathcal{H}^3}{c_s P_s - \mathcal{H}^2} \left( 1 + \frac{c_s}{c_s} \right) \left|_{k_s c_s = aH} \right.
$$

(5.16)

We have eliminated $\epsilon$ from (5.12), but two independent variables remain. To get an equation for $c_s$ we need to find $\mathcal{H}$ in terms of the observables. Since the expression for $P_h$ (5.14) only depends on $\mathcal{H}$, it can be used to find the Hubble parameter directly. In doing so, one find that

$$
\left. \mathcal{H}^2 \right|_{k_t = aH} = \frac{P_h(k_t)}{16}.
$$

(5.17)

where $P_h = A^{-1} P_h$. This gives the Hubble parameter as a function of the tensor mode wave number $k_t$. In order to find $\mathcal{H}^2$ as a function of the scalar mode wave number $k_s$, note that the relation between the wave number of tensor and scalar modes that exit the horizon at the same time is $k_t = k_s c_s$. Therefore, we can obtain $\mathcal{H}|_{k_s c_s = aH}$ by performing the substitution $k_t \rightarrow k_s c_s$:

$$
\left. \mathcal{H}^2 \right|_{k_s c_s = aH} = \frac{P_h(k_s c_s)}{16}.
$$

(5.18)
Plugging this in for $H$ into equation (5.16) we get

\[
\frac{d\mathcal{P}_h(k_t)}{d \log k_t} = -\frac{2\mathcal{P}_h(k_t)^2}{16c_s(k_s)\mathcal{P}_s(k_s) - \mathcal{P}_h(k_t)}.
\]

(5.19)

Rearranging (5.19) we find the equation for $c_s$:

\[
c_s(k_s) = \frac{1}{16} \frac{\mathcal{P}_h(k_s c_s)}{\mathcal{P}_s(k_s)} \left( 1 - \frac{2\mathcal{P}_h(k_s c_s)}{\mathcal{P}_h(k_s c_s)} \right).
\]

(5.20)

Here, a solid circle over $\mathcal{P}_h$ will denote differentiation with respect to $\log k_t$, not $\log k_s$. Once we specify what $\mathcal{P}_s(k_s)$ and $\mathcal{P}_h(k_t)$ are, we can use the above to solve for $c_s$. Even if we can only show that $c_s \neq 1$, this will be a signal for non-minimal kinetic terms. Once we use equations (5.18) and (5.20) to get $H(k)$ and $c_s(k)$ we can use (5.11) to find the relation between the wave number and time and ultimately find $H(t)$ and $c_s(t)$.

Having successfully found $H$ and $c_s$, the next step will be to use this information to find the action. The most general single scalar field Lagrangian is a multivariable function of $\phi$ and $X$. However, since we only have $H$ and $c_s$ as functions of a single independent variable (in this case $k$) the Lagrangian can only be determined as a function of $k$: $L(k)$. To turn $L(k)$ into $L(X, \phi)$, we need to find $\phi$ and $X$ as functions of $k$, invert them to get $k(\phi)$ and $k(X)$, and substitute into $L(k)$. However, there is an ambiguity in how we substitute $k$ for $\phi$ and $X$. Whenever $k$ appears in the expression for $L(k)$, we do not know whether to substitute it with $k(\phi)$, $k(X)$ or some combination of the two. The ambiguity can be partially resolved if at the onset we specify a partition of the Lagrangian into functions that depend either entirely on $\phi$ or entirely on $X$. In light of this fact we make an ansatz:

\[
L(X, \phi) = \mathcal{A}^2 M_{pl}^2 \varrho(x, \varphi) = \mathcal{A}^2 M_{pl}^2 P(g_1(x), \ldots, g_m(x), f_1(\varphi), \ldots, f_n(\varphi)),
\]

(5.21)

where $\varphi = M_{pl}^{-1} \phi$ and $x = (\mathcal{A} M_{pl})^{-2} X$. We have also defined a dimensionless Lagrangian $\varrho$ in order to keep the exposition neat and clear. Here it is assumed that
\( P(y_1, \ldots, y_m, z_1, \ldots, z_n) \) is a known function of the \( \{y_i\} \) and \( \{z_\alpha\} \), and the functions \( \{g_i\} \) and \( \{f_\alpha\} \) are not all known. We say that a Lagrangian is *partitioned* if it is written in the form given by (5.21), and the function \( P \) is referred to as the partition of the Lagrangian.

Whether a Lagrangian has a partition depends on what function \( P \) the user assumes. The user defined function \( P \) is usually chosen according to some theoretical motivation. For example, before non-canonical kinetic terms were considered, single field inflation models were almost exclusively assumed to have the form

\[
(5.22) \quad g(x, \varphi) = x - V(\varphi).
\]

Here the partition is \( P(z_1, z_2) = z_1 - z_2 \). If one assumes canonical terms, then the potential \( V(\varphi) \) can be reconstructed from the inflationary observables using methods that have been developed previously [111]. As we have learned in recent years, other types of kinetic terms are possible. For instance, in brane inflation models the Lagrangian has the form

\[
(5.23) \quad g(x\varphi) = -f^{-1}(\varphi)\sqrt{1 - 2f(\varphi)x} + f^{-1}(\varphi) - V(\varphi).
\]

In this case the partition is \( P(z_1, z_2, z_3) = -z_2^{-1}\sqrt{1 - 2z_2z_1 + z_2^{-1} - z_3} \). A reconstruction of this action would involve finding the warp factor \( f(\varphi) \) and potential \( V(\varphi) \), which we will do in section §5.3. In each of these cases the Lagrangian’s dependence on \( x \) is assumed to be known, however, we may have a model where the dependence on \( x \) is uncertain. For example, we might have a theoretical motivation for replacing \( x \) in the Lagrangian (5.22) with \( x^2 \). Thus, we could generalize (5.22) by replacing \( x \) with an unknown function: \( g(x) \), and reconstruct \( g \) from the observables to see if the result more closely matches \( g(x) = x \) or \( g(x) = x^2 \) or something else. Therefore, our procedure can be seen as a generalization of inflationary potential reconstructions.
We say that (5.21) only partially resolves the ambiguity, because it is possible in certain circumstances that a field redefinition can leave the form of (5.21) unaltered. For an example consider the Lagrangian

\begin{equation}
\rho(x, \varphi) = f(\varphi)g(x).
\end{equation}

(5.24)

Under a general field redefinition \( \varphi = h(\tilde{\varphi}) \) the Lagrangian becomes

\[ \rho(x, \varphi) = f(h(\tilde{\varphi}))g((h'(\tilde{\varphi}))^2 \tilde{x}). \]

If the function \( g \) is such that \( g(x \cdot y) = g(x) \cdot g(y) \) then

\[ \rho(x, \varphi) = f(h(\tilde{\varphi}))g((h'(\tilde{\varphi}))^2)g(\tilde{x}) = \tilde{f}(\tilde{\varphi})g(\tilde{x}), \]

where \( \tilde{f}(\tilde{\varphi}) = f(h(\tilde{\varphi}))g((h'(\tilde{\varphi}))^2) \). Thus, not all choices for the function \( P \) lead to a unique partition between the functions of \( x \) and \( \varphi \). However, while the partition is not always unique, its uniqueness cannot be determined until the functions \( g \) and \( f \) have been found. Case in point, in order to maintain the partition in our example (5.24), we needed to assume that \( g \) had the property that \( g(x \cdot y) = g(x) \cdot g(y) \). However, this assumes that we know something about the function \( g \), which would defeat the purpose of using the observables to derive \( g \) in the first place.

The first equation for the reconstructed action will be obtained from the definition of the sound speed, which is given by

\begin{equation}
c_s^2 = \frac{L_x}{L_x + 2XL_{XX}}.
\end{equation}

(5.25)

Assuming that the Lagrangian has the form (5.21), equation (5.25) can now be used to find a differential equation for the \( g_i \)'s as a function of time. After some work, this equation is given by

\begin{equation}
\langle \ddot{g}P \rangle = \frac{1}{2X} \left[ \frac{2}{x \ddot{x}} + \frac{1}{c_s^2} - 1 \right] - \langle \dot{g}P\dot{g} \rangle,
\end{equation}

(5.26)
where a dot denotes differentiation with respect to the dimensionless time $\tau = \mathcal{A}t$.

Here, we have used the following short hand notation:

\[
\langle \ddot{g} P \rangle = \sum_{i=1}^{m} \ddot{g}_i P_i, \quad \langle \dot{g} P \rangle = \sum_{i=1}^{m} \dot{g}_i P_i, \quad \langle \dot{g} P \dot{g} \rangle = \sum_{i,j=1}^{m} \dot{g}_i \dot{g}_j P_{ij},
\]

where $P_i$ and $P_{ij}$ are $P_i = \frac{\partial P}{\partial g_i}$ and $P_{ij} = \frac{\partial^2 P}{\partial g_i \partial g_j}$. Equation (5.26), however, is incomplete since $x$ is not known explicitly. To turn (5.26) into a more usable form, we need to obtain an equation for $x$ and its derivatives in terms of the known quantities $\epsilon$ and $\mathcal{H}$. To find such a formula let’s write out the expression for the energy density of a general single scalar field action (5.1):

\[
\rho = 2x L_x - L. \quad (5.27)
\]

As a consequence of the Friedmann equations, $\rho + p$ is proportional to $\frac{d\mathcal{H}}{dt}$:

\[
-\frac{d\mathcal{H}}{dt} = \frac{\rho + p}{2M_{pl}^2}. \quad (5.28)
\]

Solving for $\rho + p$ in (5.27), and substituting the result into (5.28) we find that

\[
x = -\frac{M_{pl}^2}{L_x} \frac{d\mathcal{H}}{dt} = \frac{\epsilon M_{pl}^2 \mathcal{H}^2}{L_x}, \quad (5.29)
\]

where in the last step we have used the definition of the slow roll parameter. With (5.21) as our assumed form of the Lagrangian, this algebraic equation for $x$ becomes a first order differential equation:

\[
\frac{x}{\dot{x}} = \frac{\epsilon \mathcal{H}^2}{\langle \dot{g} P \rangle}. \quad (5.30)
\]

Differentiating this equation, one can find an expression for $\ddot{x}$. After substituting the results of these relations, equation (5.26) becomes

\[
\langle \dot{g} P \rangle = \frac{2c_s^2 \mathcal{H}^2}{(1 + c_s^2)} \left( \dot{\eta} \mathcal{H} - 2c_s \mathcal{H} - \frac{1}{\langle g P \rangle} \sum_{\alpha=1}^{n} \dot{f}_\alpha \langle \dot{g} P \rangle_\alpha \right), \quad (5.31)
\]

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where $\tilde{\eta} = \frac{i}{H}$ and $\langle \dot{g} P \rangle_{\alpha}$ denotes partial differentiation of the quantity $\langle \dot{g} P \rangle$ with respect to $f_{\alpha}$. We refer to (5.31) as the sound speed reconstruction equation.

The non-gaussianity parameter $f_{NL}$ can also be used to find an equation for the functions $g_i$ and $f_{\alpha}$. Following from the ansatz of the curvature perturbation (5.6), $f_{NL}$ is determined by the behavior of the curvature three point function:

\begin{equation}
\langle \zeta(k_1)\zeta(k_2)\zeta(k_3) \rangle = -(2\pi)^7 \delta^{(3)}(k_1 + k_2 + k_3) P_s(K) \frac{3f_{NL}(K)}{10} \prod_i k_i^3,
\end{equation}

where $k_i = |k_i|$ and $K = k_1 + k_2 + k_3$. As we can see from (5.32), the $f_{NL}$ will depend on the size and shape of the triangle formed by the three scales of the three point function. In [35] the authors found an expression for $f_{NL}$ for general single field actions (5.1) when the three scales form an equilateral triangle:

\begin{equation}
f_{NL} = \frac{35}{108} \left( \frac{1}{c_s^2} - 1 \right) - \frac{5}{81} \left( \frac{1}{c_s^2} - 1 - 2\Lambda \right),
\end{equation}

where

\begin{equation}
\Lambda = \frac{X^2 L_{XX} + \frac{2}{3} X^3 L_{XXX}}{X L_X + 2X^2 L_{XX}}.
\end{equation}

To get $f_{NL}$ as a function of $K$, we have to evaluate (5.33) at the time when the scale $K$ passes outside of the sound horizon: $Kc_s = aH$. The scalar power spectrum depends on a single scale $k$, which has a one-to-one mapping with the time through the relation $kc_s = aH$. However, since $f_{NL}$ really depends on three different scales, the mapping between time and scale is not as straightforward. When the delta function in (5.32) is taken into account, the non-gaussianity still depends on three degrees of freedom: the magnitude of two of the scales and the angle between them [116]. To simplify matters, two of these three degrees of freedom will be fixed, so as to make $f_{NL}$ a univariate function. Since the equilateral configuration: $k_1 = k_2 = k_3$, has been very widely studied [35, 116], we will take $f_{NL}$ to be the non-gaussianity of
the equilateral bispectrum. The equilateral non-gaussianity will be a function of $k_{NL}$, which is the length of the sides of the equilateral triangle. Since the non-gaussianity freezes in when the scale $K$ leaves the sound horizon, the scales at which the non-gaussianity and the scalar perturbation freeze in are not the same but instead related by $3k_{NL} = k_s$. After some work, one can use (5.33) and (5.30) to show that the $g$’s and $f$’s satisfy the equation

$$\langle \dot{g}P \rangle = \frac{16\epsilon \mathcal{H}^3 c_s^2}{55} \frac{1}{1 - c_s^2 - \frac{972}{275} c_s^2 f_{NL}}$$

(5.35)

$$\times \left( \tilde{\kappa} + \frac{\epsilon c_s^2 \mathcal{H}}{\langle \dot{g}P \rangle^3} \sum_{\alpha=1}^{n} \hat{f}_{\alpha} \left[ \langle \dot{g}P \rangle \left( \langle \dot{g}P \rangle_{,\alpha} + \langle \dot{g}P \dot{g} \rangle_{,\alpha} \right) - \langle \dot{g}P \rangle_{,\alpha} \left( \langle \dot{g}P \rangle + \langle \dot{g}P \dot{g} \rangle \right) \right] \right),$$

where $\tilde{\kappa} = \frac{\dot{c}_s}{\mathcal{H}c_s} = -\frac{\kappa(1-\epsilon)}{1-\kappa}$. This is the non-gaussianity reconstruction equation. Note that (5.35) is only well defined if the last line is nonzero. If the last line does vanish and the right hand side of (5.31) is nonzero then $1 - c_s^2 - \frac{972}{275} c_s^2 f_{NL} = 0$, and the non-gaussianity (5.33) only depends on the functions $g_i$ and $f_{\alpha}$ through the sound speed $c_s$. We will discuss such a case in the next section. Finally, another relation between the $f$’s and $g$’s can be derived by combining the Friedmann equations (5.2) and (5.3):

$$P(g_1, ..., g_m, f_1, ..., f_n) = (2\epsilon - 3)\mathcal{H}^2.$$  

(5.36)

The upshot is that we now have four equations: (5.30), (5.31), (5.35), and (5.36), which when combined can be used to find $x(t)$ (and by extension $\varphi(t)$) and three of the functions $f_{\alpha}$ and $g_i$. With more observational inputs it may be possible to determine even more $f$ and $g$ functions, but for now we will be content with what we have. In what follows, we will consider different, specific scenarios for the action and show how the action in each can be determined from the observables.
5.2.1 Examples

In the case where the Lagrangian (5.21) has only one function of $x$ the equations (5.35) and (5.31) take on a much simpler form:

\begin{equation}
\dot{g}_P = \frac{2\epsilon c_s^2 \mathcal{H}^2}{(1 + c_s^2)} \left( \eta \mathcal{H} - 2\epsilon \mathcal{H} - \frac{1}{P_g} \sum_{\alpha=1}^{n} \dot{f}_\alpha P_{g\alpha} \right),
\end{equation}

\begin{equation}
\dot{g}_P = \frac{16\epsilon \mathcal{H}^3 c_s^2}{55} \left( 1 - c_s^2 - \frac{972}{275} c_s^2 f_{NL} \right) \left( \bar{\kappa} + \frac{\epsilon c_s^2 \mathcal{H}}{P_g^3} \sum_{\alpha=1}^{n} \dot{f}_\alpha \left[ P_{g\alpha} P_{g\alpha} - P_{g\alpha} P_{g\alpha} \right] \right),
\end{equation}

where $P_g = \frac{\partial P}{\partial g}$, $P_{gg} = \frac{\partial^2 P}{\partial g^2}$, $P_{g\alpha} = \frac{\partial^2 P}{\partial g \partial f_\alpha}$, and $P_{g\alpha} = \frac{\partial^3 P}{\partial g^2 \partial f_\alpha}$. As we mentioned earlier not all forms of the action will yield an equation for the functions $g_i$ and $f_\alpha$. In particular if the action is such that $P_g P_{g\alpha} = P_{g\alpha} P_{gg}$ for each $\alpha$, and if the sound speed is constant, then (5.38) is not well defined. To see why, let’s assume that $c_s$ is constant. As a result, according to the definition of the sound speed

\begin{equation}
c_s^2 = \frac{\vartheta_x}{\vartheta_x + 2x \vartheta_{xx}} \quad \Rightarrow \quad \vartheta(x, \varphi) = f_1(\varphi) x^{\frac{1+c_s^2}{2c_s^2}} + f_2(\varphi),
\end{equation}

where the $f_1$ and $f_2$ are integration constants and in general will be functions of $\varphi$ only. We already know that with this form of the action, $P_g P_{g\alpha} = P_{g\alpha} P_{gg}$. Thus the term in (5.38) inside the large parentheses vanishes, however, the right hand side of equation (5.37) does not vanish. Therefore, we expect the denominator: $1 - c_s^2 - \frac{972}{275} c_s^2 f_{NL}$, in equation (5.38) to vanish. Indeed if we use the formulas (5.33) and (5.34) for $f_{NL}$ we find that

\begin{equation}
f_{NL} = \frac{275}{972} \left( \frac{1}{c_s^2} - 1 \right).
\end{equation}

This relation between $c_s$ and $f_{NL}$ holds regardless of what the functions $f_1$ and $f_2$ in (5.39) are. It might be argued that if $c_s$ is constant then the Lagrangian (5.39) can be assumed and the remaining equations can be used to find $f_1$ and $f_2$. This however
is not the case since we have already used the sound speed equation to find \( g(x) \).

This can be confirmed if one assumes the Lagrangian (5.39). With (5.39) as our Lagrangian, equation (5.31) is equivalent to the time derivative of equation (5.30). Thus, there are really only two equations: either (5.30) or (5.31), and equation (5.36). Therefore, only one of the two \( f_1 \) and \( f_2 \) can be solved for.

There is still yet another potential complication that may arise, specifically when the Lagrangian takes the form \( \varrho(x, \varphi) = f(\varphi)g(x) \). Using (5.36) to find \( \dot{f} \) in terms of \( \dot{g} \), the equations (5.37) and (5.38) become

\[
\dot{g}f = \frac{6\epsilon\tilde{\eta}H^3c_s^2}{3(1 + c_s^2) - 2\epsilon},
\]

\[
\dot{f} = \frac{16\epsilon\tilde{\kappa}H^3c_s^2}{55} - \frac{1}{1 - c_s^2 - \frac{972}{275}c_s^2f_{NL}}.
\]

As with the previous case, (5.42) is not defined when \( \tilde{\kappa} = 0 \). Furthermore, the first equation (5.41) is also undefined when \( \tilde{\eta} = 0 \). Since this is equivalent to \( \epsilon = \text{constant} \), let’s assume that \( \epsilon \) is constant to see which type of action this corresponds to. From (5.29)

\[
\epsilon = \frac{3xL_x}{\rho} \Rightarrow \left( 2 - \frac{3}{\epsilon} \right) x\varrho_s = \varrho \Rightarrow \varrho(x, \varphi) = f_1(\varphi)x\frac{\epsilon}{2}.
\]

Notice, that this Lagrangian is a special case of the \( c_s = \text{constant} \) Lagrangian (5.39) with \( f_2 = 0 \). Thus, \( \epsilon = \text{constant} \) implies that \( c_s = \text{constant} \). However, the converse of this is not true if \( f_2 \neq 0 \). Since (5.37) is a well defined equation even with the Lagrangian (5.43), we suspect that the denominator in (5.41) vanishes. After calculating the sound speed with the Lagrangian (5.43) one finds that

\[
c_s^2 = \frac{2\epsilon - 3}{3}.
\]

So indeed, the denominator in (5.41) does vanish. Assuming that the equations (5.41) and (5.42) are well defined, consistency requires that the right hand sides of
these equations be equal, leading to the relation

\[ \frac{3\tilde{\eta}}{3(1 + c_s^2) - 2\epsilon} = \frac{8\tilde{\kappa}}{55 (1 - c_s^2) - \frac{972}{275} c_s^2 f_{NL}}. \]  

(5.44)

This is a consistency relation between \( f_{NL}, c_s \) and the slow roll parameters. Although this consistency relation only holds for models with the Lagrangian \( \varrho = f(\varphi)g(x) \), analogous consistency relations can be found for any model in question. In what follows, we will carry out the full derivation of the \( g \) and \( f \) functions for two special cases.

**Case 1:** \( \varrho(x, \varphi) = g(x) - V(\varphi) \)

Suppose the Lagrangian has the form

\[ \varrho(x, \varphi) = g(x) - V(\varphi). \]  

(5.45)

This type of Lagrangian corresponds to the standard scalar field action when \( g \) is the identity map: \( g(x) = x \). We refer to \( g(x) \) as the kinetic function. Notice that we have replaced what should be \( f_1 \) in our previous nomenclature with \( V(\varphi) \) in order to draw a clear analogy with the potential in the canonical scalar field action. With this type of action the equations (5.37) and (5.38) become

\[ \dot{g} = \frac{2\epsilon c_s^2 H^3 (\tilde{\eta} - 2\epsilon)}{1 + c_s^2}, \]  

(5.46)

\[ \dot{g} = \frac{16\epsilon\tilde{\kappa} H^3 c_s^2}{55} \frac{1 - c_s^2 - \frac{972}{275} c_s^2 f_{NL}}{1 + c_s^2 - \frac{972}{275} c_s^2 f_{NL}}. \]  

(5.47)

Here we have two expressions for the derivative of \( g(\tau) \). Consistency demands that the right hand sides of these equations be equal, thus we are lead to an analogue of the relation (5.44):

\[ \frac{\tilde{\eta} - 2\epsilon}{1 + c_s^2} = \frac{8\tilde{\kappa}}{55 (1 - c_s^2) - \frac{972}{275} c_s^2 f_{NL}}. \]  

(5.48)
Interestingly enough this is the same consistency relation found in [114] for general single field inflation models. However, while the relation in [114] was only approximate, in our case it is exact. This shows that the consistency relation of Bean et al. is exact in the case when the Lagrangian is of the form (5.45). Continuing with the derivation, the equations for $V(\tau)$ and $x(\tau)$ are given by (5.36) and (5.30), respectively. They read

\begin{align}
V(\tau) &= g(\tau) - (2\epsilon - 3)\mathcal{H}^2, \\
\frac{\dot{x}}{\mathcal{H}x} &= \frac{\dot{g}}{\epsilon\mathcal{H}^3} = \frac{2c_s^2(\tilde{\eta} - 2\epsilon)}{1 + c_s^2}.
\end{align}

With our equations in hand we are almost ready to solve them and find the action. However, we still lack knowledge about the $\mathcal{H}$ and $c_s$. In order to go further we need to look back to section §5.2 and in particular equations (5.18) and (5.20). In order for these equations to be of any use we need two observables as inputs. For these we will assume that the two inputs are the scalar and tensor contributions to the CMBR. Presently, it is believed that these spectra are near-scale invariant, and over a limited range of scales possess the forms

\begin{align}
P_s(k) &= e^{(n_s - 1)\log k/k_0}, \\
P_h(k) &= B e^{n_t \log k c_s(k)/k_0 c_{s0}}.
\end{align}

Here $c_{s0} = c_s(k_0)$, where $k_0$ is the fiducial scale at which $P_s = 1$. Note that we have assumed that the spectral indices have no running: i.e., $n_s, n_t = constant$. Admittedly, (5.50) and (5.51) are likely only approximate, as recent observations suggest [109]. In general, we expect that the spectral indices themselves have some dependence on the scale. Therefore, the exponents in (5.50) and (5.51) can be interpreted as a Taylor series expansion of the spectral indices around some scale $k = k_0$, truncated at the first order. The point $k_0$ around which we expand can
be any scale where the spectral indices and Hubble parameter have been measured. By going to second order in this expansion we can account for any running in the spectral indices. However, if we limit the range of $k$ accordingly, we can safely neglect any running and use (5.50) and (5.51). Just what range $k$ has to be limited to depends on the value of the coefficient of the second order term. Certainly, if the running is substantial this range will be very tightly constrained. Once we solve for $\varphi(k)$ and $x(k)$, the range of validly over $k$ will translate into a range for $\varphi$ and $x$ over which our results for the functions $g_i(x)$ and $f_\alpha(\varphi)$ can be trusted. Despite our assumption of constant spectra indices, some insight can be gained on the effects of running if one considers what happens to the results if $n_s$ and $n_t$ vary slightly. For example, the effect of running in the case of DBI inflation can be inferred by studying the dependence of the warp factor on the spectral indices. As discussed in the conclusion §5.4 running may result in a multi-throat DBI inflation scenario. Although running spectral indices is an interesting generalization, in order to better demonstrate the usefulness of this procedure we will stick with the simpler case of no running.

Equation (5.18) then tells us that the Hubble parameter is simply proportional to the square root of the tensor perturbation:

\[ \mathcal{H}(k) = \frac{\sqrt{\mathcal{P}_\mathcal{H}(k)}}{4} = \frac{\sqrt{B}}{4} e^{\frac{n_t}{2} \log \kappa c_s(k)/\kappa_0 c_s_0}. \] (5.52)

Defining $\mathcal{H}_0$ as $\mathcal{H}(k_0) = \mathcal{H}_0$, the constant $B$ is therefore $B = 16\mathcal{H}_0^2$. As it stands, (5.52) is not complete since we still do not have an expression for $c_s(k)$. To find $c_s$ we turn to equation (5.20), the solution of which gives us

\[ c_s = c_{s0} \left( \frac{c_{s0}}{\mathcal{H}_0^2 n_t - 2} \right)^{\frac{1}{n_t-2}} e^{-\frac{n_t-n_s+1}{n_t-2} \log k/k_0}. \]
Since we defined \( c_s \) as \( c_s(k_0) = c_{s0} \), consistency of our definition demands that

\[
c_{s0} = \mathcal{H}_0^2 \frac{n_t - 2}{n_t}.
\]

Note that if \( 0 < n_t < 2 \), \( c_{s0} \) is negative: a nonsense result. Therefore, we must restrict \( n_t \) to be either \( n_t < 0 \) or \( n_t > 2 \). With an expression for \( c_s \) in hand, \( \mathcal{H}(k) \) explicitly in terms of \( k \) is

\[
\mathcal{H}(k) = \mathcal{H}_0 e^{\frac{n_t}{n_t - 1} \log k / k_0}.
\]

Note that \( \epsilon \) and \( \kappa \) are constant in this case and are given by

\[
(5.53) \quad \begin{align*}
\epsilon &= \frac{\mathcal{H}^2}{c_s P_s} = \frac{n_t}{n_t - 2}, \\
\kappa &= -\frac{c_s}{c_s} = \frac{n_t - n_s + 1}{n_t - 1}.
\end{align*}
\]

Since \( \epsilon \) is a constant then \( \bar{\eta} = 0 \), which will simplify matters later when we try solve the reconstruction equations. Solving for \( \log k \) in (5.11), we find that the time dependence of \( \log k \) is

\[
(5.54) \quad \log k / k_0 = \frac{\omega}{\kappa} \log \left[ 1 + \epsilon \mathcal{H}_0 (\tau - \tau_0) \right],
\]

where \( \omega = \frac{\kappa}{\epsilon} \frac{1 - \epsilon}{1 - \kappa} = -2 \frac{(n_t - n_s + 1)}{n_t(n_t - 2)}. \) Therefore, the sound speed and Hubble parameter as functions of time are

\[
(5.55) \quad c_s(\tau) = c_{s0} \left[ 1 + \epsilon \mathcal{H}_0 (\tau - \tau_0) \right]^{-\omega}, \quad \mathcal{H}(\tau) = \frac{\mathcal{H}_0}{1 + \epsilon \mathcal{H}_0 (\tau - \tau_0)}.
\]

These are the expressions for the sound speed and Hubble parameter that will be used throughout this chapter. They are completely independent of the form of the action that we are solving for, and are determined only by the inflationary observables \( P_s \) and \( P_h \).
Before we go about solving the reconstruction equations we should point out that not all values of the spectral indices lead to realistic inflationary scenarios. As has been mentioned before, the sound horizon of the scalar fluctuations is not the same as the cosmological horizon. As a consequence it is now possible for the size of the sound horizon to increase as time progresses. Thus, the usual expectation that larger scales freeze in at the beginning of inflation and smaller scales freeze in at the end, is not always guaranteed to hold. Recall that the time dependence of the scale is given by equation (5.54). It follows that the sound horizon size depends on time like

\[
\text{Sound Horizon Size } \propto ( \gamma a H )^{-1} = (k/k_0)^{-1} = (1 + \epsilon H_0(\tau - \tau_0))^{-\frac{\omega}{\kappa}}.
\]

(5.56)

If \( \omega/\kappa > 0 \), the sound horizon decreases with time as is normally expected. However, if \( \omega/\kappa < 0 \), the size of the sound horizon increases during inflation, allowing modes that were previously frozen-in behind the horizon to reenter while inflation is still going on. This is a potential hazard, since if the horizon increased during inflation then widely separated regions in the visible universe never had an opportunity to get in thermal equilibrium with each other. This would make it difficult for inflation to explain the horizon problem, which is one of the reasons inflation was considered in the first place. Nevertheless, it might be possible for the horizon to increase during a small portion of the inflationary era, so long as the horizon is smaller at the end of inflation. Clearly, in our simple scenario with no spectral index running it is not possible to achieve this since (5.56) is either monotonically increasing or decreasing depending on the values of \( n_s \) and \( n_t \). One might suggest that by including running in the spectral indices the sound horizon could expand for a brief period of time. Upon inspection of (5.11), we can see that the only way the right hand side of (5.11) can be positive, and thus lead to a growing sound horizon, is if \( \kappa < 1 \) (keep in mind that \( H > 0 \) and \( \epsilon < 1 \) during inflation). Since the horizon has to eventually decrease, \( \kappa \)
must decrease at some point, and pass through the value \( \kappa = 1 \). This is problematic, since (5.11) is singular at \( \kappa = 1 \), which means that \( \log k(t) \) is not analytic there. Although it may not be impossible for this transition to occur, in this chapter we are assuming no running. Thus, the horizon will either be monotonically increasing or decreasing. Since we wish to model inflation we will assume that the horizon is increasing. Therefore, according to equation (5.56) the spectral indices must be fixed such that \( \omega \) and \( \kappa \) are either both positive or both negative.

Since we are considering only those models that allow for inflation, we need to be sure that the spectral indices are such that an inflationary phase is allowed. If we refer to the expression for the equation of state \( w \) we see that not all values of \( n_t \) are allowed if we want to have inflation:

\[
  w = \frac{p}{\rho} = \frac{L}{2XL_x - L} = -\frac{n_t - 6}{3(n_t - 2)}.
\]

Notice that so long as \( n_t < 2 \) the equation of state is always \( w < -\frac{1}{3} \), and so inflation will occur. Since \( c_s \propto \epsilon^{-1} \), in order for \( c_s \) to be interpreted as a sound speed, \( \epsilon \) must be positive. If we look back to equation (5.53) we find that not all values of \( n_t \) will result in a positive value for \( \epsilon \). Requiring that \( \epsilon > 0 \), we find that \( n_t \) must be either \( n_t < 0 \) or \( n_t > 2 \). Since we have already found that \( n_t > 2 \) would not lead to an inflationary solution, we conclude that \( n_t < 0 \). Recall that in the previous paragraph we found that the sound horizon could expand during inflation only if the spectral indices were chosen so that \( \omega/\kappa > 0 \). If one refers back to the definitions of \( \omega \) and \( \kappa \) in terms of the spectral indices, we can see that if \( n_t < 0 \) the scalar spectral index is required to be \( n_s < 2 \).

The sound speed (5.55) can tell us something about the expected range of validity of the scalar (5.50) and tensor spectra (5.51). If \( \omega < 0 \) then at some time \( \tau > \tau_0 \) the sound speed will be greater than one, signaling that fluctuations propagate at
superluminal speeds. Likewise, superluminal speeds also occur at times $\tau < \tau_0$ when $\omega > 0$. Therefore, (5.50) and (5.51) can only be considered approximations; reliable within a certain range of scales. Keeping in mind that the sound horizon needs to shrink during inflation, the wave number $k$ must respect the following bounds if the sound speed is to be less than the speed of light$^4$:

\[ \frac{k}{k_0} > (c_{s0})^\frac{1}{2\kappa}, \quad \text{for } \omega > 0, \]
\[ \frac{k}{k_0} < (c_{s0})^{\frac{1}{2\kappa}}, \quad \text{for } \omega < 0. \]

The only way (5.50) and (5.51) could be acceptable at all scales is if $\omega = 0$, in which case $c_s$ is a constant. If it turns out that the sound speed is not constant, (5.50) and (5.51) are most likely too naive. The most recent data from WMAP [109] suggests that the scalar spectral index may have a small but nonzero running, so we should not be surprised that our simple expressions for the perturbation spectra are not exactly correct. Regardless, scalar and tensor spectra with constant spectral indices are still a good approximation to the CMBR data. Our discussion will still be of relevance, as long as we keep in mind that the reconstructed actions are only approximations, valid over a limited range of scales.

We will now simplify our discussion by fixing the sound speed to a constant, which is achieved by setting $\omega = 0$. Although we will be considering only constant sound speeds, we will keep the value of $c_s$ arbitrary. This will allow us to find a more general solution to the reconstruction equations, while allowing us to study the limit $c_s \to 1$ and see whether the canonical action is recovered. One might object to this choice of $\omega$ on the grounds that if $c_s$ is constant the non-gaussianity reconstruction equation (5.47) will be ill-defined for reasons discussed in section §5.2.1. However, we counter

$^4$It has been recently proposed [31] that scalar field theories with non-minimal kinetic terms can allow for propagation of superluminal perturbations without violating causality.
that this is acceptable since we are assuming that only two functions \( g \) and \( V \) are unknown, thereby making the third reconstruction equation (5.47) unnecessary. It should be pointed out that while the two unknown functions can still be found when \( \omega = 0 \), the consistency relation (5.48) is no longer well defined. Once we substitute (5.55) for \( c_s \) and \( \mathcal{H} \) into (5.46) and solve for \( g(\tau) \) the result is

\[
g(\tau) = \frac{2\epsilon\mathcal{H}_0^2 c_{s0}^2}{1 + c_{s0}^2} \left( \frac{1}{(1 + \epsilon\mathcal{H}_0(\tau - \tau_0))^2} - 1 \right) + g_0. \tag{5.58}
\]

Using the Friedmann equation (5.49) we can now find \( V(\tau) \):

\[
V(\tau) = \frac{\mathcal{H}_0^2}{1 + c_{s0}^2} \left( \frac{3(1 + c_{s0}^2) - 2\epsilon}{(1 + \epsilon\mathcal{H}_0(\tau - \tau_0))^2} - 2\epsilon c_{s0}^2 \right) + g_0. \tag{5.59}
\]

Similarly, the equation for \( \dot{x} \) is given by

\[
\frac{\dot{x}}{x} = -\frac{4\epsilon c_{s0}^2}{1 + c_{s0}^2} \frac{\mathcal{H}_0}{[1 + \epsilon\mathcal{H}_0(\tau - \tau_0)]},
\]

and the exact solution for \( x \) is

\[
x(\tau) = x_0 \left[ 1 + \epsilon\mathcal{H}(\tau - \tau_0) \right]^{\frac{1 - c_{s0}^2}{1 + c_{s0}^2}}. \tag{5.60}
\]

Integrating (5.60) to find \( \varphi(\tau) \)

\[
\varphi = \frac{\dot{\varphi}_0}{\epsilon\mathcal{H}_0} \left( 1 + c_{s0}^2 \right) \left[ (1 + \epsilon\mathcal{H}_0(\tau - \tau_0))^{\frac{1 - c_{s0}^2}{1 + c_{s0}^2}} - 1 \right] + \varphi_0,
\]

where \( \dot{\varphi}_0 = \sqrt{2x_0} \). Now that we have \( x \) and \( \varphi \) as functions of time, we can invert these and substitute the results into (5.58) and (5.59) to find \( g(x) \) and \( V(\varphi) \). After carrying this out, we can combine the \( g(x) \) and \( V(\varphi) \) to arrive at the full Lagrangian:

\[
g(x, \varphi) = \frac{2\epsilon\mathcal{H}_0^2 c_{s0}^2}{1 + c_{s0}^2} \left( x/x_0 \right)^{\frac{1 - c_{s0}^2}{2c_{s0}^2}} - \frac{\mathcal{H}_0^2}{1 + c_{s0}^2} \left[ 3(1 + c_{s0}^2) - 2\epsilon \right] \left[ 1 + \epsilon\mathcal{H}_0 \frac{1 - c_{s0}^2}{1 + c_{s0}^2} \varphi - \varphi_0 \right]^{\frac{2(1 + c_{s0}^2)}{1 - c_{s0}^2}}.
\]

Here is the complete Lagrangian in the case when the sound speed is constant. Note that the final result does not depend on the integration constant \( g_0 \). This is a result
of the fact that the right hand side of equation (5.36) is independent of the initial values of the kinetic and potential functions. The only undetermined constants are the initial values of the scalar field and its derivative, and due to the attractor nature of inflation, their exact values are unimportant. Despite what was said earlier in regards to the indefiniteness of the non-gaussianity reconstruction equation (5.47), this Lagrangian does have a definite non-gaussianity given by the result in equation (5.40). It is worth noting that in the exceptional case where $c_s \to 1$:

$$\varrho(x, \varphi) = \epsilon \mathcal{H}_0^2 (x/x_0) - \mathcal{H}_0^2 (3 - \epsilon) e^{-\frac{2\epsilon \mathcal{H}_0}{\varphi_0} (\varphi - \varphi_0)},$$

we recover the canonical inflation Lagrangian with an exponential potential. If we require that the wave function retain the standard normalization then $x_0 = \epsilon \mathcal{H}_0^2$ and the Lagrangian becomes

$$\varrho(\varphi, x) = x - \mathcal{H}_0^2 (3 - \epsilon) e^{-\sqrt{2\epsilon} (\varphi - \varphi_0)},$$

(5.61)

which is the Lagrangian of power-law inflation [117]. This is a reassuring result; it confirms that in the appropriate limit, we can recover the standard inflationary action.

**Case 2:** $\varrho(x, \varphi) = f(\varphi)g(x) - V(\varphi)$

Let us now take the complexity of the action one step further and assume that there are now three unknown functions: $g$, $f$ and $V$. We define $\varrho$ as

$$\varrho(x, \varphi) = f(\varphi)g(x) - V(\varphi).$$

(5.62)

In equation (5.37) the only nonzero $P_{g\alpha}$ is the one corresponding to $f$. Thus (5.37) reduces to

$$\dot{g}f = \frac{2\epsilon c_s^2 \mathcal{H}^2}{(1 + c_s^2)} \left( \tilde{\eta} \mathcal{H} - 2\epsilon \mathcal{H} - \frac{\dot{f}}{f} \right).$$

(5.63)
Furthermore, with this Lagrangian the terms with the $f_\alpha$’s in (5.38) all vanish. The final result is simply

\begin{equation}
\frac{\dot{g}f}{\mathcal{H}} = \frac{16\epsilon\bar{\kappa}\mathcal{H}^3c_s^2}{55} \frac{1}{1 - c_s^2 - \frac{972}{275} c_s^2 f_{NL}}.
\end{equation}

Combining equations (5.63) and (5.64), $g$ decouples and we get an equation just for $f$:

\begin{equation}
\frac{\dot{f}}{\mathcal{H}f} = \tilde{\eta} - 2\epsilon - \frac{8\bar{\kappa}}{55} \frac{1 + c_s^2}{1 - c_s^2 - \frac{972}{275} c_s^2 f_{NL}}.
\end{equation}

Once we have solved for $f$ here we can substitute the solution into equation (5.64) and solve for $g$. With the solutions for these two, $V$ is found using the Friedmann equation (5.36). The final step is to find $\varphi(\tau)$ and $x(\tau)$ by solving (5.30):

\begin{equation}
\dot{x} = \frac{16\epsilon\bar{\kappa}\mathcal{H}c_s^2}{55} \frac{x}{1 - c_s^2 - \frac{972}{275} c_s^2 f_{NL}}.
\end{equation}

Let’s assume that $\bar{\kappa} \neq 0$, so that the reconstruction equations (5.64) (5.65) and (5.66) are well defined. We will again assume that the scalar and tensor perturbation spectra are given by (5.50) and (5.51). Therefore, $\mathcal{H}$ and $c_s$ are the same as those that we found earlier (5.55). However, now that we are using the non-gaussianity reconstruction equation we need to specify $f_{NL}$. In this example we will take $f_{NL} = 0$ to simplify the analysis. With these as our inputs, the reconstruction equations become

\begin{align*}
\frac{\dot{f}}{\mathcal{H}f} &= -2\epsilon + \frac{8\epsilon \omega}{55} \frac{1 + c_s^2}{1 - c_s^2}, \\
\frac{\dot{g}f}{\mathcal{H}} &= -\frac{16\epsilon^2 \omega \ c_s^2 \mathcal{H}^3}{55} \frac{1}{1 - c_s^2}, \\
\frac{\dot{x}}{\mathcal{H}x} &= -\frac{16\epsilon \omega \ c_s^2}{55} \frac{1}{1 - c_s^2}.
\end{align*}
Here, we have used the fact that \( \tilde{\eta} = 0 \) and \( \tilde{\kappa} = -\epsilon \omega \). Each of these has an analytic solution. They are

\[
\begin{align*}
  f(\tau) &= f_0 \left[ \frac{c_s^2(\tau)}{c_s^2(0)} \right]^{\frac{1}{2}} \frac{4}{\pi} \left( \frac{1 - c_s^2(\tau)}{1 - c_s^2(0)} \right)^{8/55}, \\
  g(\tau) &= g_0 + \frac{2\mathcal{H}_0(1 - c_s^{-2})^\frac{8}{55}}{c_s f_0} \left[ \frac{c_s^2(\tau)}{c_s^2(0)} \right]^{-\frac{1}{55}} F(c_s^{-2}(\tau)) - F(c_s^{-2}(0)), \\
  x(\tau) &= x_0 \left( \frac{1 - c_s^2(\tau)}{1 - c_s^2(0)} \right)^{\frac{8}{55}}.
\end{align*}
\]

(5.68) \quad (5.69) \quad (5.70)

Here we have defined \( F(x) \) as

\[
F(x) = _2F_1 \left( \frac{4}{55}, \frac{63}{55}, \frac{59}{55}; x \right),
\]

where \( _2F_1 \) is a hypergeometric function. Note that we can find a complete expression for \( g(x) \). We simply have to solve for \( c_s(\tau) \) in (5.70) to get \( c_s^2(x) \), which is

\[
c_s^2(x) = 1 - (1 - c_s^2(0)) \left( \frac{x}{x_0} \right)^{-55/8},
\]

(5.71)

and then substitute this for \( c_s(\tau) \) in (5.69) to get \( g(x) \). Interestingly enough, \( g(x) \) is independent of \( \omega \), so taking the \( \omega \to 0 \) limit here is trivial. Since the solution for \( g(x) \) is in terms of hypergeometric functions, to get a better idea of what \( g(x) \) looks like we expand around \( c_s = 1 \), and thus obtain

\[
g(x) = g_0 + \frac{\epsilon \mathcal{H}_0^2}{f_0 c_s} \frac{x - x_0}{x_0} + \frac{\epsilon \mathcal{H}_0^2(1 - c_s)}{2585 f_0 c_s} \left[ 2145 - 2209 \frac{x}{x_0} + 64 \left( \frac{x_0}{x} \right)^{\frac{47}{2}} \right] + \mathcal{O}((1 - c_s)^2).
\]

(5.72)

Let’s take a moment to comment on the analytic behavior of \( g(x) \). In fig. 5.1 the exact functional behavior of \( g(x) \) is shown along with the approximate expression (5.72). As fig. 5.1 and the approximation (5.72) suggest, the behavior of \( g \) is nearly linear with respect to \( x \), when \( c_s \) is close to one. However, for reasons that will be clear shortly, the limit \( c_s \to 1 \) does not necessarily mean that the Lagrangian will be
Figure 5.1: Plot depicting the function $g(x)$. This plot was made with $c_{s0} = 0.5$, $\epsilon = 0.1$, $f_0 = 1$, $g_0 = 0$ and $\phi_0 = 1$. The exact behavior of $g(x)$ is contrasted against the approximation (5.72). The behavior of $g(x)$ is very linear except for small deviations for $x < x_0$. Note that at $x_0(1 - c_{s0}^2)^{8/55} \approx 0.48$ the plot of the exact behavior of $g(x)$ stops abruptly as a result of the fact that $g$ becomes non-real in this region.

linear in $x$. Another interesting feature of $g(x)$ is that it becomes non-real for values of $x$ less than $x_0(1 - c_{s0}^2)^{8/55}$. This implies a lower bound on the values of $x$, which is a behavior that is observed in the solution (5.70). This lower bound is a result of the fact that for $x < x_0(1 - c_{s0}^2)^{8/55}$ the sound speed squared would be negative according to (5.71).

As for the functions $f(\phi)$ and $V(\phi)$, one cannot find analytic expressions for these like we did for $g(x)$. Once we integrate $x(\tau)$ to find $\phi(\tau)$, we can see why:

\begin{equation}
\varphi = \varphi_0 + \frac{\sqrt{2} x_0 c_{s0} (1 - c_{s0}^2)^{4/55}}{H_0^3 (1 + \frac{8 \omega}{55})} \left[ \left[ \frac{c_{s0}^2 (\tau)}{c_{s0}^2} \right]^{-\frac{1}{55} - \frac{4}{55}} F_\omega(c_{s0}^{-2}(\tau)) - F_\omega(c_{s0}^{-2}) \right],
\end{equation}

where we have again shortened things by defining $F_\omega$ as

\[ F_\omega(x) = 2 F_1 \left( \frac{4}{55} + \frac{1}{2 \omega}, \frac{4}{55}, \frac{59}{55} + \frac{1}{2 \omega}; x \right). \]

Since $\varphi$ is such a complicated function there is no way to invert (5.73) to get time as an analytic function of $\varphi$. Therefore, we are forced to either evaluate $f(\varphi)$ and $V(\varphi)$...
numerically, or make an approximation for $\tau(\varphi)$. Since we will be interested in finding a correspondence with the example in the previous section, we will approximate $\varphi(\tau)$ in the $\omega \ll 1$ limit. The result of this approximation is

\begin{equation}
\varphi(\tau) = \varphi_0 + \dot{\varphi}_0 \tau - \frac{8\dot{\varphi}_0 \omega}{55\epsilon H_0} \frac{c^2_{s0}}{1 - c^2_{s0}} [(1 + \epsilon H_0 \tau) \log(1 + \epsilon H_0 \tau) - \epsilon H_0 \tau] + O(\omega^2),
\end{equation}

where we have set $\tau_0 = 0$. To get $\tau(\varphi)$ we will drop all $\omega$ dependent terms from (5.74), so that $\varphi(\tau)$ is a linear function of $\tau$. This approximation turns out to be remarkably accurate even at late times, since the higher order terms in (5.74) scale only logarithmically with $\tau$. Now that we have at least an approximate expression for $\tau(\varphi)$, $f(\varphi)$ can be found by replacing $c_s(\tau)$ with

\begin{equation}
c_s(\varphi) = c_{s0} \left(1 + \frac{\epsilon H_0}{\varphi_0} (\varphi - \varphi_0)\right)^{-\omega}.
\end{equation}

The exact behavior of $f(\varphi)$ was evaluated numerically and the results are shown in fig. 5.2. Since we will be taking the $\omega \to 0$ limit later, we make a further approximation of $f(\varphi)$ by Taylor expanding around $\omega = 0$:

\begin{equation}
f(\varphi) \approx \frac{f_0}{(1 + \frac{\epsilon H_0}{\varphi_0} (\varphi - \varphi_0))^2} + \frac{8\omega f_0 (1 + \frac{\epsilon H_0}{\varphi_0} (\varphi - \varphi_0))}{55} \frac{1 + \frac{\epsilon H_0}{\varphi_0} (\varphi - \varphi_0)}{1 - c^2_{s0}} \log(1 + \frac{\epsilon H_0}{\varphi_0} (\varphi - \varphi_0)).
\end{equation}

Notice that the second term diverges at $c_{s0} = 1$. Therefore, although the higher order terms in (5.72) vanish when $c_{s0} = 1$, when the limit $c_{s0} \to 1$ is taken the product $fg$ will retain the nonlinear $x$ terms. This is why the Lagrangian may not be linear in $x$ even when $c_{s0} = 1$. To find the potential we use the Friedmann equation (5.36). Doing so requires us to find $g$ as a function of $\varphi$, which we find by replacing $c_s(\tau)$ in (5.69) with $c_s(\varphi)$ (5.75). The potential, Taylor expanded around $\omega = 0$, is

\begin{equation}
V(\varphi) \approx \frac{g_0 f_0 - (2\epsilon - 3)H^2_0}{(1 + \frac{\epsilon H_0}{\varphi_0} (\varphi - \varphi_0))^2} + \frac{8\omega g_0 f_0 (1 + \frac{\epsilon H_0}{\varphi_0} (\varphi - \varphi_0))}{55} \frac{1 + \frac{\epsilon H_0}{\varphi_0} (\varphi - \varphi_0)}{1 - c^2_{s0}} \log(1 + \frac{\epsilon H_0}{\varphi_0} (\varphi - \varphi_0)).
\end{equation}
Figure 5.2: Plot depicting the function $f(\phi)$ for different values of $\omega$. This plot was made with $c_s=0.5$, $\epsilon=0.1$, $\varphi_0=0$, $\dot{\varphi}_0=1.0$, $f_0=1.0$ and $g_0=0$. We have also included a plot of the function $f(\phi)$ in (5.80) for comparison.

Note that the individual functions $g$, $f$ and $V$ depend on the arbitrary integration constants $f_0$ and $g_0$, even though the action that is composed of them does not. If we are interested in just finding the action, fixing $f_0$ and $g_0$ would be a moot point. However, it does raise the matter of how one separates the action into kinetic and potential terms. For example, suppose we separate the kinetic function into a constant and a “$x$-dependent” piece:

$$g(x) = c + G(x).$$

The constant $c$ is arbitrary and can be adjusted to any given value by absorbing the difference into $G(x)$. Substituting the right hand side of (5.78) for $g(x)$, the Lagrangian (5.62) becomes

$$g(x, \varphi) = f(\varphi)G(x) - V(\varphi) + cf(\varphi).$$

With the Lagrangian written in this way, it would make more sense to define $G(x)$
as the kinetic function and define the potential as

\begin{equation}
(5.79) \\
v(\phi) = V(\phi) - cf(\phi).
\end{equation}

In the case where \( f \) is constant (such as the example in the previous section) then the above redefinition only amounts to a uniform shift in the potential. However, if \( f \) is non-constant then the behavior of the potential can change drastically. Although none of the CMBR data are sensitive to changes in \( c \), it is possible to find a value for \( c \) by requiring that in the appropriate limit, the action becomes equivalent to the canonical action. We will define this as the canonical limit of the action. Before we determine \( c \) by this method we need to confirm that the Lagrangian (5.62) is canonically equivalent to the canonical Lagrangian (5.61) when the sound speed is constant and equal to one.

If we turn our attention back to our approximations for \( f(\phi) \) and \( V(\phi) \), we notice that taking \( c_{s0} = 1 \) leads to divergent results. These divergences are understandable since the reconstruction equations (5.67) are divergent when \( c_s = 1 \). However, if we set \( \omega = 0 \) in (5.76) and (5.77), it’s possible to take \( c_{s0} = 1 \) and still obtain a well defined result. Doing so results in the following for the functions \( g, f \) and \( V \):

\begin{equation}
(5.80) \\
g(x) \approx g_0 + \frac{\epsilon \mathcal{H}_0^2 x - x_0}{f_0} x_0, \quad f(\phi) \approx \frac{f_0}{(1 + \frac{\epsilon \mathcal{H}_0}{\varphi_0}(\phi - \varphi_0))^2}, \quad V(\phi) \approx \frac{f_0 g_0 - (2\epsilon - 3)\mathcal{H}_0^2}{(1 + \frac{\epsilon \mathcal{H}_0}{\varphi_0}(\phi - \varphi_0))^2},
\end{equation}

where now \( \mathcal{H}_0^2 = \epsilon \) since \( c_{s0} = 1 \). We refer the reader to fig. 5.2 for a comparison of \( f(\phi) \) in (5.80) and \( f(\phi) \) for general values of \( \omega \) and \( c_{s0} \). The Lagrangian in the \( c_{s0} \rightarrow 1 \) limit when \( \omega = 0 \) is

\begin{equation}
(5.81) \\
g(x, \phi) = \frac{\epsilon \mathcal{H}_0^2}{(1 + \frac{\epsilon \mathcal{H}_0}{\varphi_0}(\phi - \varphi_0))^2} x_0 - \frac{(3 - \epsilon)\mathcal{H}_0^2}{(1 + \frac{\epsilon \mathcal{H}_0}{\varphi_0}(\phi - \varphi_0))^2}.
\end{equation}
This Lagrangian can be related to the canonical Lagrangian (5.61) through the field redefinition defined by

\[ 1 + \frac{\epsilon H_0}{\varphi_0} (\varphi - \varphi_0) = e^{\sqrt{2} (\tilde{\varphi} - \tilde{\varphi}_0)}. \]

Under this redefinition, the new Lagrangian is

\[ \varrho(\tilde{x}, \tilde{\varphi}) = \tilde{x} - (3 - \epsilon) H_0^2 e^{-\sqrt{2} (\tilde{\varphi} - \tilde{\varphi}_0)}. \]

This is the same as the canonical Lagrangian (5.61) found in the first example. The ability to redefine the field so as to obtain the canonical action was only possible in the limit where the sound speed is equal to one and constant. In the case where the sound speed is constant but not equal to one, the actions are not canonically equivalent. To see why this is, note that a necessary condition for two actions to be canonically equivalent is that they both lead to the same observables. Since the non-gaussianity in the first case is given by (5.40), whereas the non-gaussianity in the second case is assumed to vanish, we can see that the two actions are not canonically equivalent except when \( c_{s0} = 1 \), which is the canonical limit.

To ensure a smooth transition to the canonical action we must separate the kinetic function as we did in (5.78) so that \( G(0) = 0 \) in the \( \omega \to 0 \) and \( c_{s0} \to 1 \) limits. Upon inspection of \( g(x) \) in (5.80) we see that the redefined kinetic function \( G(x) \) is

\[ G(x) = g(x) - \frac{f_0 g_0 - \epsilon H_0^2}{f_0}. \]

In doing so the potential is redefined according to (5.79) as

\( (5.82) \quad v(\varphi) = V(\varphi) - \frac{g_0 f_0 - \epsilon H_0^2}{f_0} f(\varphi). \)

There is a subtlety in this analysis that should be addressed. In order to reclaim the canonical Lagrangian we needed to take the limits \( \omega \to 0 \) and \( c_{s0} \to 1 \) simultaneously.
In our case we took that limit by setting $\omega = 0$ and then letting $c_{s_0}$ approach one. However, this is by no means the only way to take the limit. For example, we could have approached the limit by setting $\omega = 1 - c_{s_0}^2$ and then take the limit as $c_{s_0}$ goes to one. Had we taken the limit from a different direction it is possible that the Lagrangian that resulted could have been different from the canonical Lagrangian (5.61). After some inspection, it can be shown that under a field redefinition $\varphi = h(\tilde{\varphi})$ such that $f^{-1}(\varphi) = (h'(\tilde{\varphi}))^2$, the potential in the canonical limit is given by

$$V(\tilde{\varphi}) = (3 - \epsilon)\mathcal{H}_0^2 e^{-\frac{2\mathcal{H}_0}{\epsilon_0}(\tilde{\varphi} - \tilde{\varphi}_0)} + (g_0 f_0 - \epsilon \mathcal{H}_0^2) \frac{f(\varphi)}{f_0}.$$  

Here $f(\tilde{\varphi}) = f(h(\tilde{\varphi}))$ is the function $f$ when the canonical limit is taken. It is simple to show that the canonical limit of $f$ is not unique, which means that the potential is also not unique. However, if we redefine our potential according to (5.82) instead, the new potential $v$ is unique, and the Lagrangian that results is canonically equivalent to (5.61).

**Case 3: $f_1(\varphi)$ and $f_2(\varphi)$ Unknown**

We now bring up a case that will be of particular interest to reconstructions of the DBI action. We start by assuming that the Lagrangian $\varrho(x, \varphi)$ has the form

$$\varrho(x, \varphi) = P(x, f_1(\varphi), f_2(\varphi)), $$

where $f_1$ and $f_2$ are unknown functions of $\varphi$. Unlike the previous cases, the functional dependence of the Lagrangian with respect to $x$ is assumed to be known exactly. In this case it is possible to obtain a set of algebraic equations of the two unknowns $f_1$ and $f_2$. The first of these equations can be most easily obtained by going back to the original definition of the sound speed (5.25):

$$c_s^2 = \frac{L_x}{L_x + 2x L_{xx}} \Rightarrow P_{xx} = \frac{1}{2x} \left( \frac{1}{c_s^2} - 1 \right) P_x.$$  

(5.83)
This equation together with (5.30) and the Friedmann equation (5.36) are enough to find \( f_1(\varphi) \) and \( f_2(\varphi) \) in terms of the observables. In the next section we will see explicitly how the equations (5.83), (5.30) and (5.36) come together to reconstruct the DBI action from the power spectrum data.

5.3 DBI inflation

In realistic string and M-theories, the number of space-time dimensions is 10 or 11 dimensions. The extra 6 or 7 dimensions are compactified to small sizes, leaving the effective theory at low energies a theory of physics in four dimensions. The various moduli that control the shape (complex structure moduli) and size (Kähler moduli) of the internal space, also determine the nature of the four-dimensional low-energy effective theory. Therefore, fixing these moduli is an important step in establishing a connection between string theory and the standard model. In recent years, much attention has been paid to flux compactifications as a potential means of stabilizing string moduli\(^5\). In a flux compactification, various fluxes wrap around closed cycles in the internal manifold creating a potential for the complex structure moduli. The best known of these takes place in type IIB string theory. Here the internal space is six-dimensional Calabi-Yau and the 3-form fluxes \( F_3 \) and \( H_3 \) create a superpotential that fixes the complex structure [119]. These 3-form fluxes source a warping of the geometry of the internal manifold. In the type IIB flux compactification, the ansatz of the line element is taken as

\[
d s_{10}^2 = h^{-1/2}(y)g_{\mu\nu}dx^\mu dx^\nu + h^{1/2}(y)g_{mn}dy^m dy^n. \tag{5.84}
\]

Here \( h \) is the warp factor which is sourced by the fluxes and varies only along the dimensions of the internal manifold. In DBI inflation, which we will be considering

\(^5\)For a review of flux compactifications see [118].
in this section, the local geometry of the internal manifold is a Klebanov-Strassler throat geometry [120], and is described by the metric

\[ g_{mn}dy^m dy^n = dr^2 + r^2 ds^2_{X_5}. \]

Here \( ds^2_{X_5} \) is the line element of a five-dimensional manifold \( X_5 \), which forms the base of the KS throat. The coordinate \( r \) runs along the depth of the throat. For our purposes we will only consider motion along \( r \) and integrate over the base manifold \( X_5 \). The warping of the internal space creates a natural realization of the Randall-Sundrum model [123], and has also provided model builders with a new approach to developing string theory based models of inflation [124]. The most popular inflation model that makes use of this warping is DBI inflation [27, 28], which is the primary focus of this section.

In the simplest DBI inflation models a \( D3 \) brane travels along the \( r \) direction, either into or out of the KS throat. The \( D3 \) brane extends into the three non-compact space dimensions and is point-like in the internal manifold. The standard DBI action for the \( D3 \) brane is

\[
S_{DBI} = -\int d^4x\sqrt{-g} \left[ f^{-1}(\phi)\sqrt{1 - 2f(\phi)X - f^{-1}(\phi) + V(\phi)} \right].
\]

(5.85)

Here \( \phi = \sqrt{T_3}r \) (where \( T_3 \) is the \( D3 \) brane tension) is a rescaling of the coordinate \( r \) and will play the role of the inflaton. The quantity \( f^{-1} = T_3h^{-1} \) is the rescaled warp factor. The metric \( g_{\mu\nu} \) that appears in (5.85) is the metric on the 3 + 1 dimensional non-compact subspace which describes the geometry of our familiar 4 dimensional space-time. We will continue to assume that the geometry of the 3 + 1 dimensional subspace is described by the FRW metric with zero curvature. The energy density

\[ \rho = \frac{\dot{\phi}^2}{2} + V(\phi). \]

(5.86)

Fluctuations of the brane position along the transverse directions of the KS throat have been mentioned as a possible source of entropy perturbations [121, 122]. These could serve as a further constraint on the form of the action.
and pressure in the non-compact subspace due to the brane are given by

\begin{equation}
\rho = f^{-1} (\gamma - 1) + V,
\end{equation}

\begin{equation}
p = (\gamma f)^{-1} (\gamma - 1) - V.
\end{equation}

Here $\gamma$ is a new parameter, not found in the standard canonical inflation. In terms of the quantities in the DBI action, $\gamma$ is defined as

\begin{equation}
\gamma = \frac{1}{\sqrt{1 - 2f(\phi)X}}.
\end{equation}

The $\gamma$ defined here is analogous to the Lorentz factor in special relativity, and will hence-forth be referred to as the Lorentz factor. The Lorentz factor places an upper limit on the speed of the brane as it travels through the KS throat. Since the speed of the brane is limited, this allows one to get a sufficient amount of inflation even with potentials that would be considered too steep to use in standard canonical inflation.

In our study of the DBI model we will be assuming that the scalar and tensor spectra are approximately (5.50) and (5.51), respectively. With these as our inflationary observables, we found that $\epsilon$ was constant (5.53). The fact that $\epsilon$ is a constant indicates that inflation will not end on its own, and instead some other mechanism such as $D3\overline{D3}$ annihilation [125] must be used to provide a graceful exit. Since our study is concerned more with the physics during inflation, this topic will not be addressed further. We will now present a generalized DBI action, and show how it is reconstructed from the inflationary observables.

### 5.3.1 A Generalized DBI Model

Having sketched out the general method for reconstructing different types of inflationary actions in section §5.2, it is now time to apply these methods to a DBI-type
Lagrangian given by

\[
\varrho(x, \varphi) = P(x, \mathcal{F}(\varphi), \mathcal{V}(\varphi)),
\]

where the partition \( P \) is given by

\[
P(z_1, z_2, z_3) = -z_2^{-1} \left( \sqrt{1 - 2z_2z_1} - 1 \right) - z_3.
\]

Here \( \mathcal{F}(\varphi) = \mathcal{A}^2 M_{pl}^2 f(\varphi) \) is the (dimensionless) warp factor in the throat, and \( \mathcal{V}(\varphi) = (\mathcal{A} M_{pl})^{-2} V(\varphi) \) is the (dimensionless) potential. In the KS throat geometry the warp factor is taken to be \( \mathcal{F} \propto \varphi^{-4} \). The potential \( \mathcal{V} \) is assumed by many to be quadratic in \( \varphi \). For the purposes of this study we will not assume a priori any form for the functions \( \mathcal{F} \) and \( \mathcal{V} \), and instead allow the inflationary observables to determine them.

Now that we have established the general form of the action, we can use the procedure outlined in section §5.2.1 to find \( \mathcal{F} \) and \( \mathcal{V} \). Turning to equations (5.83) and (5.30) we find the relations

\[
\frac{L_{xx}}{L_x} = \frac{\mathcal{F}}{1 - 2\mathcal{F} x} = \frac{1}{2x} \left( \frac{1}{c_s^2} - 1 \right), \quad x = \epsilon \mathcal{H}^2 \sqrt{1 - 2\mathcal{F} x}.
\]

Solving for \( \mathcal{F} \) and \( x \):

\[
\mathcal{F}(\varphi) = \frac{1 - c_s^2}{2\epsilon \mathcal{H}^2 c_s}, \quad x = \epsilon \mathcal{H}^2 c_s.
\]

Comparing the second equation above with (5.30) and recalling the definition of (5.87), we find that

\[
c_s = \frac{1}{\gamma}.
\]

This result is characteristic of DBI inflation and holds regardless of the warp factor and potential used. Having found \( \mathcal{F} \), equation (5.36) tells us what \( \mathcal{V} \) is:

\[
\mathcal{V}(\varphi) = 3\mathcal{H}^2 + \frac{1}{\mathcal{F}} \left( \frac{1}{c_s} - 1 \right).
\]
Having already found the expression for $F$ in (5.90), we can now write down the important reconstruction equations for $V$ and $F$ in terms of $H, c_s$ and $\epsilon$:

\begin{align}
V(k) &= H^2 \left(3 - \frac{2\epsilon}{1 + c_s}\right), \\
F(k) &= \frac{1 - c_s^2}{2\epsilon c_s H^2}.
\end{align}

To turn $F$ and $V$ into functions of $\varphi$ we need to integrate our solution for $x$ (5.90) to find $\varphi$. Taking equation (5.30) to find an expression for $\dot{\varphi}$, we find that in the case of DBI inflation

\begin{equation}
\dot{\varphi} = \pm \sqrt{2\epsilon c_s H}.
\end{equation}

The sign of the right hand side of the equation is ambiguous, due to the square root taken to get this equation from (5.30). The sign is left arbitrary for now and will be specified later based on the requirement that $\varphi$ be positive. Once we solve for $\varphi(\tau)$ in (5.94) and invert to get $\tau(\varphi)$, we can then find a solution for $V(\varphi)$ and $F(\varphi)$. Now that we have laid the ground work for generating the functions of the generalized DBI action, the next section will show how the perturbation spectra are used to obtain explicit expressions for $F(\varphi)$ and $V(\varphi)$.

5.3.2 The Warp Factor and Potential in DBI Inflation

In this section we will now use the program that was laid out at the end of the previous section to find an exact solution for the warp factor and potential in the Lagrangian (5.88). We will again assume that the scalar and tensor spectra have a power-law dependence with respect to the scale $k$. Therefore, the sound speed, Hubble parameter and $\epsilon$ are the same as those found in section §5.2.1. Thus, the potential as a function of time is

\begin{equation}
V = \frac{H_0^2}{(1 + \epsilon H_0(\tau - \tau_0))^2} \left(3 - \frac{2\epsilon}{1 + \frac{H_0^2}{\epsilon}(1 + \epsilon H_0(\tau - \tau_0))^{-\omega}}\right),
\end{equation}
and the warp factor as a function of time is

\[ F = \frac{(1 + \epsilon H_0(\tau - \tau_0))^{\omega + 2}}{2H_0^4} \left( 1 - \frac{H_0^4}{\epsilon^2} (1 + \epsilon H_0(\tau - \tau_0))^{-2\omega} \right). \]

After substituting our solutions for $H(\tau)$ and $c_s(\tau)$ into the equation of motion for $\varphi(\tau)$ we get

\[ \dot{\varphi} = \pm \sqrt{2} H_0^2 (1 + \epsilon H_0(\tau - \tau_0))^{-\frac{3}{2}} - \frac{\omega}{2}. \]

Once we integrate this expression we can obtain an answer for $\varphi(\tau)$. As a matter of convenience we will set the value of the resulting integration constant to zero. Later, once we have found $F$ and $V$, we will see that this choice allows for a correspondence between the reconstructed functions and their theoretically derived counterparts. After integrating (5.97) we find that

\[ \varphi(\tau) = \int^\tau \dot{\varphi} d\tau = \pm \frac{2\sqrt{2} H_0}{\epsilon \omega} (1 + \epsilon H_0(\tau - \tau_0))^{-\omega/2}. \]

Since we are interested in eventually connecting the reconstructed action with the standard DBI model we need to keep the inflaton, which is just a rescaled radial coordinate, positive. The sign that we choose in (5.98) will therefore depend on the sign of $\omega$. We can write the general solution as

\[ \varphi(\tau) = \varphi_0 (1 + \epsilon H_0(\tau - \tau_0))^{-\omega/2}, \]

where

\[ \varphi_0 = 2\sqrt{2} \left| \frac{H_0}{\epsilon \omega} \right|. \]

In the case where $\omega > 0$ the field $\varphi$ decreases monotonically to zero as time passes, which implies that the brane is falling into the throat. This corresponds to the UV DBI scenario. If on the other $\omega < 0$, then $\varphi$ increases monotonically with time and
the brane falls out of the throat, which corresponds to IR DBI inflation. Solving for

\[ 1 + \epsilon H_0 (\tau - \tau_0) = \left( \frac{\varphi}{\varphi_0} \right)^{-2} . \]

Substituting this for \( 1 + \epsilon H_0 (\tau - \tau_0) \) in the expressions we found for the potential

and warp factor we find that \( V \) as a function of \( \varphi \) is

\[ V = H_0^2 \left( \frac{\varphi}{\varphi_0} \right)^{\frac{4}{3}} \left( 3 - \frac{2\epsilon}{1 + \frac{H_0^2}{\epsilon}(\frac{\varphi}{\varphi_0})^{2}} \right), \]

and the warp factor as a function of \( \varphi \) is

\[ \mathcal{F} = \frac{1}{2H_0^4} \left( \frac{\varphi}{\varphi_0} \right)^{-2+\frac{4}{3}} \left[ 1 - \frac{H_0^4}{\epsilon^2} \left( \frac{\varphi}{\varphi_0} \right)^4 \right]. \]

Furthermore, when \( \gamma \) is expressed as a function of \( \varphi \), it takes on a very simple form:

\[ \gamma = \gamma_0 \left( \frac{\varphi_0}{\varphi} \right)^2 = \frac{8}{\epsilon \omega^2} \left( \frac{1}{\varphi_0} \right)^2,\]

where \( \gamma_0 = \frac{1}{c_{so}} \). It is interesting to note that (5.100) is the same as the approximate

results found in the theoretically inspired DBI model [27]. The potential and warp

factor derived here are the same as those found in [126]. There the authors recon-

structed the potential and warp factor by assuming that the equation of state \( w = \frac{p}{\rho} \)

was a constant and that \( \varphi \propto \tau^{-\omega/2} \). In contrast, we have reconstructed the potential

and warp factor under the assumption that the scalar and tensor perturbations are

(5.50) and (5.51). The non-gaussianity in this DBI model is the same result that one

comes across in the literature [35]:

\[ f_{NL} = \frac{35}{108} \left( \frac{1}{c_s^2} - 1 \right). \]

This particularly simple result is a general feature of DBI inflation, and is indepen-

dent of the warp factor and potential. This result for the non-gaussianity (5.101)
also follows from consistency of the reconstruction equations (5.31) and (5.35). Thus, (5.101) can be viewed as a consistency relation analogous to those found in (5.44) and (5.48). An interesting generalization to consider is

\[ \varrho(x, \varphi) = P(g(x), F(\varphi), V(\varphi)), \]

where \( P(z_1, z_2, z_3) \) is defined by (5.89). As a consistency check, one can easily show that if the non-gaussianity is equal to (5.101), and the potential and warp factor are given by (5.95) and (5.96), respectively, then \( g(x) = x \) is a solution to the reconstruction equations (5.31) and (5.35).

Having found the potential and warp factor as functions of the inflaton, we can now say that our task is at an end. Amazingly enough, despite the complicated form of the reconstruction equations an exact solution for \( V \) and \( F \) was available even for semi-realistic scalar and tensor spectra. In the next section we will discuss the properties of the reconstructed action, and its correspondence with the theoretically derived action of DBI inflation.

5.3.3 Discussion

In section §5.2.1 we found that not all values of the spectral indices lead to inflationary and/or physically sensible actions. Specifically, we showed that unless \( n_t < 0 \) the matter described by the action was unable to drive an inflationary phase. Furthermore, when this constraint on \( n_t \) was considered in conjunction with the requirement that the sound horizon decrease as inflation occurs, the scalar spectral index had to be bounded like \( n_s < 2 \). These results hold for any reconstructed action that was derived assuming the observational inputs (5.50) and (5.51). However, even if these constraints are satisfied, it is not guaranteed that the reconstructed action is physically sensible when interpreted in the context of a given theoretical
construction. For example, if we are to interpret the action reconstructed in this
section as a DBI action, then $\gamma > 0$. Doing so would be contrary to its definition
(5.87) within the context of DBI inflation. In this case, since $\gamma = \frac{1}{c_s}$ the fact that
we already have enforced the constraint $c_s > 0$ automatically keeps $\gamma$ positive. We
will see later, however, that the constraints found in §5.2.1 are not sufficient for our
reconstructed action to be interpreted as a DBI action.

First, let’s consider the warp factor (5.99) and what constraints it places on the
observables. Suppose $n_t < n_s - 1$. In this case the exponent of $\varphi$ in front of the square
brackets in (5.99) will always be negative. Therefore, for small values of $\varphi$ the leading
order behavior of $\mathcal{F}$ will go like $\varphi^{-a}$ where $a > 0$. Thus, the warp factor increases as
we fall into the throat, which is what we would expect for a warped compactification
in string theory. On the other hand, if $n_t > n_s - 1$ the leading order behavior of $\mathcal{F}$
will be $\varphi$ to some positive or negative power, depending on the relative difference
between $n_s$ and $n_t$. If the difference between $n_t$ and $n_s$ is too small, then to leading
order, $\mathcal{F}$ will scale like $\varphi$ to some positive power. This indicates that the warp factor
gets smaller as we reach the bottom of the throat, which is a scenario that is difficult
to embed into a string theory compactification. However, if the difference between
$n_s$ and $n_t$ is large enough, then it is possible to get a more sensible solution where
$\mathcal{F} \sim \varphi^{-a}$. In general, the condition that $\mathcal{F}$ increases as we approach the bottom of
the throat implies that

$$-1 - \frac{2}{\omega} < 0 \quad \Rightarrow \quad \omega < -2 \quad \text{or} \quad \omega > 0.$$  \hspace{1cm} (5.102)

The regions in the $n_s$-$n_t$ parameter space where the condition (5.102) is satisfied
are shown in fig. 5.3. The region shown in light grey in fig. 5.3 is defined by
$\omega > 0$, or equivalently $n_t < n_s - 1$. This region corresponds to the UV phase of DBI
inflation. The region in dark grey is defined by $\omega < -2$, or equivalently
$n_t > -\frac{n_s - 1}{n_s - 3}$.
and corresponds to the IR phase. As fig. 5.3 illustrates, if \( n_s \) is restricted to the presently favored value \( n_s \sim 0.96 \), then the value of \( n_t \) is tightly constrained in the IR region but relatively unrestricted in the UV phase. Furthermore, since \( n_t < 0 \), only those models in the UV phase can have a blue-tilt. Recent CMBR data favors a blue-tilted spectrum, but only if there is a running spectral index \([109]\). Although running spectral indices would be an interesting extension of this analysis, we will leave this topic to future studies.

An unpleasant feature of the warp factor (5.99) is that it becomes negative when \( \varphi > \varphi_0 H_0 / \sqrt{\epsilon} \). This is particularly distasteful since the metric (5.84) depends on \( f^{1/2} \), which means that at sufficiently large \( \varphi \) the metric is imaginary. The values of \( \tau \) where the warp factor is positive are given by

\[
\tau - \tau_0 > \frac{c_s^{1/2} - 1}{\epsilon H_0} \quad \text{for } \omega > 0 \quad \text{(UV)},
\]

\[
\tau - \tau_0 < \frac{c_s^{1/2} - 1}{\epsilon H_0} \quad \text{for } \omega < 0 \quad \text{(IR)}.
\]
One can show using equation (5.54) that this is equivalent to the bounds in (5.57).
This is no coincidence; it is a result of the fact that $F$ is proportional to $1 - c_s^2$ (5.90). The absence of superluminal propagation is equivalent to requiring that the warp factor is positive. Even though causality can still be preserved in the case of superluminal propagation [31], if we want to interpret the action in question as a model of DBI inflation, then $c_s < 1$ in order for this interpretation to be consistent.

Therefore, the bounds (5.103) suggest that the perturbation spectra (5.50) and (5.51) can only be used as approximations over a limited range of scales. The inequalities in (5.103) tell us that $F$ is a valid warp factor towards the end of inflation in the case of UV DBI, and at the beginning of inflation in IR DBI. The time $\tau_0$ at which the initial conditions are specified should be at the beginning of inflation in the UV scenario, and at the end in the case of IR DBI. If we choose $\tau_0$ in this manner then the approximations for the perturbation spectra (5.50) and (5.51) will lead to a realistic warp factor for the entire duration of the inflationary episode.

It is worth asking if the $F$ that we have derived in (5.99) can approximate the AdS warp factor derived from theory. It is clear from (5.99) that this can be achieved if and only if $\omega = 2$. However, the only way we can get $\omega = 2$ is if either i.) $n_t = 1$ ii.) $n_t \to \infty$ or iii.) $n_s = 1$. As we have already seen, case i.) is unphysical, and case ii.) is difficult to imagine taking place. While case iii.) is unlikely to be true exactly, it is nevertheless the more realistic of the three, especially when you consider that observation suggests that $n_s \approx 0.96$. If we do set $n_s = 1$, the warp factor and potential become

$$F(\varphi) = \frac{2}{\epsilon^4} \frac{1}{\varphi^4} - \frac{1}{2\epsilon^2},$$

$$V = \frac{3\epsilon^2}{2} \varphi^2 - \frac{\epsilon^3 \varphi^2}{1 + \frac{\epsilon}{2} \varphi^2}.$$
In the case where the field range is small these are approximately

\begin{equation}
(5.104)
\begin{align*}
f(\phi) & \approx \frac{2M_{\text{pl}}^2}{A^2 \epsilon^4 \phi^4}, \\
V(\phi) & \approx \frac{\epsilon^2 A^2 (3 - 2\epsilon)}{2} \phi^2,
\end{align*}
\end{equation}

where we have reverted back to the standard, dimensionful $f$, $V$ and $\phi$ for clarity’s sake. It is a bit of a surprise that in the process of trying to recover the AdS warp factor we have stumbled upon the commonly used potential in UV DBI inflation. If we take (5.104) and demand that it is consistent with the theoretical result we can arrive at a condition on $\epsilon$ in terms of the D3 brane charge. Recall that in the KS throat $f(\phi)$ is given by

\[ f(\phi) = \frac{2T_3 R^4}{\phi^4}, \]

where

\[ T_3 = \frac{1}{(2\pi)^3 g_s (\alpha')^2}, \]

and

\[ R^4 = 4\pi g_s N(\alpha')^2 \frac{\pi^3}{\text{Vol}(X_5)}. \]

Consistency with (5.104) demands that

\[ \frac{2M_{\text{pl}}^2}{A^2 \epsilon^4} = \frac{\pi N}{\text{Vol}(X_5)} \Rightarrow \epsilon \approx \frac{10^2}{N^{1/4}}. \]

In order to get an inflationary phase $N \approx 10^{10}$, putting us well within the range of validity for the supergravity approximation. While it is interesting that the standard D3 brane DBI model can be recovered from a near-scale invariant scalar power spectrum, it has been acknowledged that this inflation model is problematic. In [127] Baumann and McAllister found that while present bounds on non-gaussianity imply that $N \lesssim 38$, primordial perturbations imply that $N \gtrsim 10^8 \text{Vol}(X_5)$. These two limits
are incompatible unless \( \text{Vol}(X_5) \lesssim 10^{-7} \). It is not clear that such a space could be naturally embedded into a string theory compactification. More general warp factors and potentials have been considered in [113]. There it was found that models could not simultaneously satisfy bounds on the field range and observational bounds on the non-gaussianity. Therefore, even though our warp factor and potential matches the theoretically based predictions, the problems inherent in the DBI model carry over into its generalizations.

5.4 Conclusion

In this chapter we have presented a method for deriving the actions of single field inflation models using CMBR data. This method allows one to derive up to three unknown functions of the action using the scalar perturbation \( P_s \), tensor perturbation \( P_t \) and the non-gaussianity \( f_{NL} \). After stating the reconstruction equations, we carried out the reconstruction procedure for two simple examples. For the purposes of the reconstruction, we assumed that the scalar and tensor spectra were power-law dependent on the scale \( k \), with the spectral indices kept as free parameters. In the first example we assumed that the Lagrangian had the form shown in equation (5.45), and used the reconstruction equations to obtain the action as a function of the spectral indices. In this example there were only two unknown functions, thus the reconstruction equations also led to a consistency relation (5.48) between the \( f_{NL} \), \( c_s \) and the slow roll parameters. However, this consistency relation is only well defined when the sound speed is not a constant.

In the second example, the action depended on three unknown functions and therefore required all three reconstruction equations. In order to simplify the discussion we took as our input for the non-gaussianity \( f_{NL} = 0 \). Although we were
unable to express the action in terms of elementary functions we were able to obtain
the action numerically and approximately assuming $c_s 0 \approx 1$ and $\omega \approx 0$. We showed
that in the limit where $c_s$ is constant and equal to one, the action in this example
was canonically equivalent to the canonical action derived at the end of the previous
section. In discussing this example we also pointed out possible ambiguities in the
program relating to how one defines a separation between the kinetic and potential
terms.

In section §5.3, we used the procedure to derive and study the warp factor and
potential in a generalized DBI inflation model. Again, we assumed that both of the
perturbation spectra scaled like $k$ to some power. Exact expressions for the warp
factor and potential were then derived, each having an explicit dependence on the
spectral indices. The demand for a physically sensible DBI inflation model placed
constraints on the spectral indices. In addition we found that the derived action
approximates the original UV DBI inflation model in the case where $\phi \ll M_{pl}$ and
$n_s = 1$. Unfortunately, the problems that have plagued UV DBI inflation are still
present in our case.

This procedure was shown to be useful in studying how the action of a general
inflation model depends on the observables. For example, we found that if the scalar
and tensor perturbation spectra went like $k$ to an arbitrary power, the reconstruction
would lead to a realistic inflationary model only if $n_t < 0$ and $n_s < 2$. Furthermore, to
keep the speed of fluctuations from becoming superluminal, the range of $k$ over which
the approximations for the spectra (5.50) and (5.51) are taken, had to be limited.
When we reconstructed a generalized DBI action in section §5.3, further constraints
were needed to keep the action compatible with an interpretation of DBI inflation.
Specifically, we found that in order to keep the warp factor positive, the field range
had to be limited. Furthermore, in the theoretically motivated DBI model, the warp factor increases as we reach the bottom of the warped throat. In order for this to be true in our reconstruction, the spectral indices needed to satisfy the additional constraints: $\omega < -2$ or $\omega > 0$.

In this chapter we have only considered the simplest of the DBI inflation models, which unfortunately suffers from several inconsistencies. However, there are many extensions of the D3 brane DBI model that can circumvent some of the problems of the original. Some of these extensions include using wrapped D5 branes [128, 129], multiple D3 branes [130], and multiple throats [57]. Each of these models has its potential advantages and drawbacks. Applying our reconstruction procedure may help to further elucidate their relative strengths and weaknesses. Furthermore, we have limited ourselves to perturbations with simple power-law behavior. However, this naive assumption may be incorrect. It is easy to imagine that the spectral indices themselves are also scale dependent. Based on the results of this chapter we can predict what kind of effect a running spectral index would have on the physics of the underlying models. For instance, in the generalized DBI model it is possible for the spectral indices to change during inflation in such a way as to pass from the IR to the UV phase. Transition between phases would correspond to a completely different physical scenario, one where the brane falls out of one throat and back into another. Therefore, running spectral indices would describe multi-throat DBI inflation. A model which has so far been shown to be internally consistent [131].

This study has also raised some other questions that may be worth investigation. In particular what is the relation between actions that yield the same observables.

It may be possible to define a group of transformations which leave the perturbation

\footnote{Fig. 5.3 implies that inflation can change between UV and IR phases only if it passes through the exactly scale-invariant point: $(n_s, n_t) = (1, 0)$.}
spectra and the non-gaussianity invariant. Such a set of transformations would allow us to classify actions based on the observables they yield. Another interesting possibility that came out of this study is the idea of using the reconstruction equations as a way to generate consistency relations between $f_{NL}$, the sound speed $c_s$ and the slow roll parameters. These questions will be left for future studies.
6.1 Introduction

The observation of the accelerated expansion of the universe [9] has been one of the most important recent discoveries in cosmology. Many possible explanations have been put forward, which may be classified under two general classes: models with a cosmological constant; or dynamical models of dark energy. Their common feature is that they provide fluids with negative pressure to drive the acceleration. Among the dynamical dark energy models, only $k$-essence [19] has the advantage of explaining not only the current phase of accelerated expansion, but also the coincidence problem; i.e., why the cross-over from the matter dominated era to the current era happened so recently in the past. This explanation, however, is not without its own problems, as was first pointed out in [30]. It was shown there that in order to solve the coincidence problem, the universe had to go through an era where the speed of propagation of the k-essence fluctuations must become superluminal. This problem was addressed in [31], where it was shown that at the classical level, superluminal propagation does not necessarily imply causal paradoxes. In particular, propagation in a k-essence background does not have any additional causal difficulties over general relativity, where the only problems are associated with space-times
that admit closed time-like curves. In the course of this analysis, the authors found a very interesting way of describing the propagation of k-essence fluctuations in terms of an emergent metric that depends not only on the space-time metric, but also the background cosmological k-essence scalar. Thus, one may think of this non-trivial, Lorentz violating cosmological k-essence background as the “aether” in which matter perturbations propagate. This emergent metric description is used in this chapter to couple the k-essence background to neutrinos.

If there is a scalar field pervading the universe, then the effective field theory viewpoint implies that it must undergo interactions with the matter that is present. The question of the observability of dark energy directly through its couplings to ordinary matter is an important one [132, 133, 134, 135, 136, 137, 138, 139, 140, 141, 142, 143]. This chapter attempts to address aspects of it in the context of k-essence. On short distance scales the universe is inhomogeneous with, in particular, plenty of black holes. Thus, the interactions of dark energy with black holes could be one way to study the above mentioned question. In this chapter, however, we will be concerned only with the effect of the k-essence aether on the propagation of neutrinos. The coupling of fermions to this background is in itself an interesting question from the theoretical point of view. Most studies of the interaction of dark energy with fermions couple them through a Yukawa-like interaction [142], which is quite reasonable. However, in this chapter, we do this differently using the vierbeins constructed out of the emergent metric. Throughout this chapter the k-essence field is treated strictly as a background; however, we see no reason not to treat it as a dynamical field. As we argue in section §6.2, within the effective field theory methodology, the terms with higher derivatives of the fermion fields do not give rise to ghosts. The main focus of this chapter is on looking for observable consequences of dark energy,
so in this case the fermion in question is the neutrino, which we show undergoes
flavor oscillations when traveling through the k-essence aether. However, the way we introduce the fermion/k-essence coupling could be used to obtain new types of interactions between dark energy and other forms of matter including dark matter. The emergent metric from [31], which is used throughout this chapter to couple the dark energy to neutrinos, is covariantly constant. In a future publication we will discuss how to consider even more generalized couplings by introducing torsion in the emergent space-time.

It has been noticed previously [31], in the context of the propagation of k-essence fluctuations in a classical background that defines the emergent space-time, that Lorentz invariance is lost when the speed of propagation of the fluctuations (the speed of sound $c_s$), is different from the speed of light. The same is true for the propagation of neutrinos in this background. In fact, all of the physically interesting results that we obtain in this chapter are present only when $c_s \neq c$; i.e., when there is Lorentz violation. Non-trivial neutrino flavor oscillations require, in addition, non-diagonal flavor couplings of neutrinos to the k-essence background. In the past, various models of neutrino oscillations have been considered that require an explicit violation of Lorentz invariance [52], or the equivalence principle [49, 50]. The energy dependence of the oscillation length is the same in these models as the one considered in this chapter. In this sense our model may be considered a theoretically and phenomenologically motivated manifestation of the same phenomenon. We should emphasize that we do not have any violation of the equivalence principle; the emergent metric which contains contributions from the k-essence background can be different for different flavors of neutrinos. This is made possible in our model by a flavor non-diagonal coupling in the part of the emergent metric involving the
k-essence background only.

K-essence is a theory with non-canonical kinetic terms, and coupling it to neutrinos through the vierbein of the emergent metric alters the speed of propagation of the neutrinos. The consequent dispersion relations are analyzed in section §6.2. The data from supernova 1987a [144, 145, 146] is then used as an input to constrain some of the parameters of our model. In section §6.3 we consider the possibility of neutrino oscillations induced by their coupling to the k-essence background. The more interesting and novel case is when the different neutrino flavors couple with different strengths to the k-essence field. Oscillations are induced essentially due to the fact that the speeds of propagation of the different neutrino species in the background aether are consequently different. In this case neutrinos would oscillate even if they were massless. In section §6.3.2 we discuss this case in some detail, and obtain a general formula for the oscillation probability with massive neutrinos. Our results are quantitatively different than the case of flavor oscillations with only massive neutrinos. In particular, the oscillation length varies with the inverse power of the neutrino energy. Such a behavior is ruled out by the data from the Kamiokande experiment [147]. Thus, we are able to place bounds on the allowed strengths of the k-essence coupling to neutrinos. In particular, the data strongly favor diagonal flavor couplings of k-essence to neutrinos. As a preliminary to this analysis, in section §6.3.1 we discuss the case of neutrino oscillations with massive neutrinos, but with equal coupling strengths of all flavors to the k-essence background. Here we find a rather simple modification of the well known formula for the flavor oscillations of massive neutrinos; with the only difference arising due to the fact that the neutrinos travel along geodesics in the emergent space-time. In section §6.4 we present our conclusions. Certain technical details of the coupling of Dirac fermions to the
emergent metric are relegated to an appendix.

6.2 Neutrino Coupling To a K-essence Background

Before we can talk about neutrino interactions with k-essence, we should first review the latter [19, 31]. In general k-essence is a theory of a scalar field with non-canonical kinetic terms. The Lagrangian of a single k-essence scalar \( \phi \) is usually denoted by a single function \( L(X, \phi) \), where \( X = \frac{1}{2} \nabla_\mu \phi \nabla^\mu \phi \). For a given solution \( \phi \) to the k-essence equations of motion, the behavior of perturbations \( \pi = \delta \phi \) in the k-essence field around the background \( \phi \) can be described by a canonical scalar field action, but with the space-time metric \( g_{\mu\nu} \) replaced by an emergent space-time metric \( G_{\mu\nu} \) given by [31]

\[
G_{\mu\nu} = \Omega^{-2} \left( g_{\mu\nu} + \frac{\delta^2}{2X} \nabla_\mu \phi \nabla_\nu \phi \right),
\]

where \( \Omega^2 = \frac{\delta}{L_X} \), and \( \delta^2 = \frac{L_X}{L_X + 2XL_{XX}} \). In the case where \( X > 0 \), the parameter \( \delta \) is equal to the sound speed of k-essence fluctuations \( c_s \), which is defined as

\[
(6.1) \quad c_s^2 = \left( \frac{\partial p_k}{\partial \rho_k} \right)_{\phi}.
\]

Here \( \rho_k \) and \( p_k \) denote the energy density and pressure of the k-essence background, and the subscript \( \phi \) signals that (6.1) should be evaluated while holding \( \phi \) constant. In the case where \( \nabla_\mu \phi \) is a time-like vector (i.e., \( X > 0 \)), it can be shown that \( p_k = L \) and \( \rho_k = 2XL_X - L \). Therefore, it follows that

\[
c_s^2 = \delta^2 = \frac{L_X}{L_X + 2XL_{XX}}.
\]

Since \( g_{\mu\nu} \) is replaced by \( G_{\mu\nu} \) in the \( \pi \) field action, the characteristics of \( \pi \) follow the geodesics of \( G_{\mu\nu} \) and not \( g_{\mu\nu} \). This interesting fact implies that \( \pi \) has a different

\[1\] Derivatives of \( L \) with respect to \( X \) are denoted by a subscript, so that \( L_X = \frac{\partial L}{\partial X} \) and \( L_{XX} = \frac{\partial^2 L}{\partial X^2} \).
causal structure than all other fields. In particular, if in some frame $g_{\mu\nu} = \eta_{\mu\nu}$ and the background k-essence field is uniform, then the emergent space-time metric is given by

$$G_{\mu\nu} dx^\mu dx^\nu \propto c_s^2 dt^2 - d\mathbf{x}^2.$$  \hfill (6.2)

The metric $G_{\mu\nu}$ defines a different causal structure than $g_{\mu\nu}$. The “light” cones of $G_{\mu\nu}$ are defined by characteristics that have velocity $c_s$ instead of the speed of light. It was shown in [31] that even if $c_s$ exceeded the speed of light, causality in a k-essence theory would still be preserved despite the superluminal speed of k-essence fluctuations. This fact can be roughly understood by thinking of the k-essence metric (6.2) as the space-time interval in special relativity but with a different value for the speed of light. Therefore, causality is preserved in k-essence for much the same reason that it is preserved in special relativity. However, in order for this rational to hold the k-essence action must satisfy certain constraints. These constraints arise from the need to make $G_{\mu\nu}$ have the proper signature. This translates to the requirement that the k-essence action satisfies

$$1 + \frac{2X L_{XX}}{L_X} > 0.$$  \hfill (6.3)

Note that this condition is equivalent to the stability constraint: $c_s^2 > 0$.

In this chapter we will take inspiration from the k-essence perturbation action, and consider the possibility of other fields coupling to $G_{\mu\nu}$. In particular, we will take the action of a neutrino coupled to a gravitational metric $g_{\mu\nu}$ and replace this with the k-essence metric $G_{\mu\nu}$. The action of our hypothetical k-essence coupled (Dirac) neutrino $\nu$ is given by

$$S = \int d^4x E\nu [i\tilde{\gamma}^\mu D_\mu - M] \nu,$$  \hfill (6.4)
where $E = \det E^a_\mu$ and $\tilde{\gamma}^\mu = E^\mu_a \gamma^a$, and $\gamma^a$ are the standard gamma matrices. The vierbein field $E^\mu_a$ of the emergent space-time geometry and its inverse $E^a_\mu$ are given by

\begin{equation}
E^\mu_a = \Omega \left( e^\mu_a + \frac{g u}{2X} e^\rho \phi \nabla^\mu \phi \right), \quad E^a_\mu = \Omega^{-1} \left( e^a_\mu + \frac{g \tilde{u}}{2X} e^\rho \phi \nabla^\mu \phi \right).
\end{equation}

Here we have defined $u = \frac{1}{3} - 1$, and $\tilde{u} = -\frac{u}{1+gu}$. Further note that $\nabla^\mu \phi = g^{\mu\nu} \nabla_\nu \phi$. Also, we have included a coupling constant $g$, which accounts for the interaction strength of the k-essence background to the neutrinos. The emergent metric $G_{\mu\nu}$ is still covariantly constant as in [31].

At this point we would like to emphasize two features of the action. First, we note that when $c_s \neq 1$ the model is not invariant under Lorentz transformations [31]. As we will see in detail in section §6.3, all the physical effects that we discuss in this chapter are consequences of this Lorentz violation in the sense that they vanish at $c_s = 1$. Other features, like non-diagonal flavor couplings, are also important to get non-trivial flavor oscillations; however, Lorentz violations must always be present. Secondly, in this model we treat the k-essence background as a classical field, which does not experience any appreciable back-reaction from the neutrino field. This allows us to treat the neutrino/k-essence coupling as a contribution to the kinetic term in the neutrino action. If the k-essence field were dynamical, this would lead to higher order derivatives of the neutrino field in the k-essence equation of motion, which could potentially create ghosts in the quantum theory. This is not a problem in this chapter since we treat the k-essence scalar strictly as a classical background. However, we would like to emphasize that within the effective field theory methodology, treating the k-essence scalar as a dynamical field would not give rise to such problems in any case. The terms with higher derivatives of the fermion field would be considered as higher order in the low-energy effective action.
expansion.

The derivative operator $\mathcal{D}_\mu$ in (6.4) represents the spinor covariant derivative with respect to the emergent k-essence background described by $E_\mu^a$. The proper definition of $\mathcal{D}_\mu$ and its specific form in the case of a general k-essence field in a flat space-time background are given in appendix §B. From here on out we will ignore the spinor connection term in $\mathcal{D}_\mu$, which we justify on the basis that higher derivatives of the k-essence field are negligibly small at the present time in most models of k-essence. By definition, the emergent space-time metric is given by $G_{\mu\nu} = E_\mu^a E_\nu^b \eta_{ab}$, which is

\begin{equation}
G_{\mu\nu} = E_\mu^a E_\nu^b \eta_{ab} = \Omega^{-2} \left( g_{\mu\nu} + \frac{\ddot{w}}{2X} \nabla_\mu \phi \nabla_\nu \phi \right),
\end{equation}

where $\ddot{w} = 2g\ddot{u} + g^2\dddot{u}^2$. Note that if $g = 1$, then $\ddot{w} = \delta^2 - 1$, which is the standard result for the emergent metric in k-essence theories. The inverse of $G_{\mu\nu}$ is given by

\begin{equation}
G^{\mu\nu} = \Omega^2 \left( g^{\mu\nu} + \frac{w}{2X} \nabla^\mu \phi \nabla^\nu \phi \right),
\end{equation}

where $w = 2gu + g^2u^2 = -\frac{\ddot{w}}{1+\ddot{w}}$. The determinant of $G_{\mu\nu}$ is given by

\begin{equation}
E^2 = -\det G_{\mu\nu} = -(\det g_{\mu\nu}) \frac{\Omega^{-8}}{1+w}.
\end{equation}

In order for this modified k-essence induced metric to have the proper signature: $1+w > 0$. Note that this condition is equivalent to (6.3) in the case when $g = 1$.

The k-essence coupled Dirac equation reads

\begin{equation}
(i\gamma^\mu \mathcal{D}_\mu - M) \nu(t,x) = 0.
\end{equation}

If we square this equation we can obtain a Klein-Gordon equation for $\nu(t,x)$:

\begin{equation}
(G^{\mu\nu} \nabla_\mu \nabla_\nu + R/4 + M^2)\nu(t,x) = 0.
\end{equation}
Here, $R$ is the scalar curvature of the vierbein $E^a_{\mu}$, and is defined in terms of the connection form by

$$R = E^a_{\mu} E^\nu_b \left[ \partial_\mu \Omega^a_{\nu} - \partial_\nu \Omega^a_{\mu} + \Omega^a_{c\mu} \Omega^c_{\nu} - \Omega^a_{c\nu} \Omega^c_{\mu} \right].$$

It is evident from (6.7) that the curvature scalar acts as a mass term. However, as we discussed earlier, in most k-essence models higher derivatives of the field $\phi$ will be negligible. Therefore, we can ignore $R$ from here on out. By taking the plane wave approximation for the neutrino field, the phase of $\nu$ is proportional to $e^{-i \int p_\mu dx^\nu}$. If we assume that the interaction of the neutrino field is weak and the background geometry is flat, then the dominant space-time dependence of the neutrino field comes from the phase factor. Thus, the Klein-Gordon equation in momentum-space leads to the dispersion relation

$$G^{\mu\nu} p_\mu p_\nu - M^2 = 0.$$

Define an effective momentum $\tilde{p}_\mu$ as

$$\tilde{p}_\mu = \Omega^{-1} e^a_\mu E^\nu_a p_\nu = p_\mu + \frac{g^{\mu}}{2X} (p_\nu \nabla^\nu \phi) \nabla_\mu \phi.$$

The covariant and contravariant effective momenta are defined with respect to the space-time metric and not the emergent k-essence metric $G_{\mu\nu}$. Thus, the index on $\tilde{p}_\mu$ is raised and lowered using the space-time metric $g_{\mu\nu}$. Because of this property, it follows that

$$\Omega^2 \tilde{p}_\mu \tilde{p}^\mu = \Omega^2 g^{\mu\nu} \tilde{p}_\mu \tilde{p}_\nu = g^{\mu\nu} (e^a_\mu E^\nu_a p_\rho) (e^b_\nu E^\lambda_b p_\lambda) = \eta^{ab} E^a_\mu E^\nu_b p_\mu p_\nu = G^{\mu\nu} p_\mu p_\nu.$$

For the purposes of this chapter we will assume that the background space-time is flat so that $g_{\mu\nu} = \eta_{\mu\nu}$ and $e^a_\mu = \delta^a_\mu$. This is in fact a good approximation cosmologically for the applications we have in mind. Therefore, the dispersion relation in terms of
the effective momentum is

\begin{equation}
\tilde{p}_0 \tilde{p}^0 + \tilde{p}_i \tilde{p}^i = m^2,
\end{equation}

where \( m^2 = \Omega^{-2} M^2 \) is an effective mass that we have defined here for convenience.

We wish to use the on-shell condition above to find the particle velocity \( v \) of the neutrinos, which is represented by the group velocity \( v = \frac{\partial E}{\partial |p|} \).

Throughout this chapter we will define \( p_\mu = (E, -\mathbf{p}), \ p = |\mathbf{p}|, \ \mathbf{p} = \mathbf{p} \hat{n}, \) and \( \dot{\phi} = \nabla_0 \phi \). In the next few sections we will find the neutrino velocity for two special cases: a uniform k-essence field, and a static k-essence field, after which we will derive the velocity assuming the most general k-essence field configuration.

### 6.2.1 Simple Case: \( \phi \) is Uniform

Before we try and find the expression for the neutrino velocity with the most general k-essence field, let’s find the velocity when \( \phi \) is uniform. This is probably the most relevant case since most k-essence theories, in particular those that attempt to address the cosmological constant problem, assume that the spatial derivatives of the k-essence field are negligible compared to its time derivative [30, 19, 21, 148]. If \( \phi \) is uniform then \( \nabla_i \phi = 0 \), and thus

\[ \tilde{p}_0 \tilde{p}^0 + \tilde{p}_i \tilde{p}^i = m^2 \Rightarrow E = \frac{\sqrt{p^2 + m^2}}{1 + \frac{q \dot{\phi}^2}{2X}}. \]

By definition \( X = \frac{1}{2} \dot{\phi}^2 \). Therefore,

\begin{equation}
\frac{v_p}{p} = \frac{E}{p} = c_\nu \sqrt{1 + \frac{m^2}{p^2}},
\end{equation}

and the group velocity, which represents the neutrino particle velocity is

\begin{equation}
\frac{v_g}{p} = \frac{\partial E}{\partial p} = \frac{c_\nu}{\sqrt{1 + \frac{m^2}{p^2}}},
\end{equation}
where \( c_\nu = \frac{1}{1+gu} = \frac{c_s}{(1-g)c_s+g} \). The speed \( c_\nu \) plays the same role for massless neutrinos that the sound speed \( c_s \) does for massless k-essence perturbations. Both \( c_\nu \) and \( c_s \) represent limiting speeds that neutrinos and k-essence perturbations, respectively, are required to not exceed as measured from the frame in which the k-essence background is uniform. The fact that in general \( c_\nu \neq c_s \) is due to our inclusion of an arbitrary coupling parameter \( g \). Phenomenologically, \( u \) is small, so we easily see that if \( c_s > 1 \), then so is \( c_\nu \).

If \( m^2 = 0 \) then the neutrino velocity (6.10) will be equal to \( c_\nu \). As a massless particle coupled to the emergent background geometry, the neutrino will travel on the null geodesics of \( G_{\mu\nu} \) not \( g_{\mu\nu} \). In the uniform case, the emergent metric (6.6) with arbitrary \( g \) is

\[
G_{\mu\nu} dx^\mu dx^\nu \propto c_\nu^2 dt^2 - d\mathbf{x}^2.
\]

As one can see here, null lines in \( G_{\mu\nu} \) travel at a speed \( c_\nu \). In this sense the emergent geometry of a uniform k-essence field acts just as a minkowski space-time except that the limiting speed is now \( c_\nu \) instead of the speed of light.

### 6.2.2 Slightly Less Simple Case: \( \phi \) is Static

In direct contrast to the last case, let’s consider what happens to the neutrino velocity when \( \phi \) is time independent, but has nonzero spatial gradients. In this case \( \nabla_\mu \phi \) is a space-like vector (i.e., \( X < 0 \)) and the k-essence energy density and pressure are instead given by \( \rho_k = -L \) and \( p_k = L - 2XLX \). It follows that \( c_s^2 = \frac{LX+2XLXX}{LX} \), which means that the definition of the sound speed for a space-like k-essence field is the inverse of the sound speed for a time-like k-essence field. If we assume that there is a static, but spatially varying k-essence field, then the on-shell condition states
\[ \tilde{p}_0 p^0 + \tilde{p}_i p^i = m^2 \quad \Rightarrow \quad E^2 = p^2 \left[ 1 + w \cos^2 \theta \right] + m^2, \]

where \( \theta \) is the angle between \( p \) and \( \nabla \phi \). Note that in this case \( X = -\frac{1}{2} |\nabla \phi|^2 \). The neutrino phase velocity is then given by

\[ v_p = \frac{E}{p} = \sqrt{1 + w \cos^2 \theta + \frac{m^2}{p^2}}, \]

and the group velocity is

(6.11) \[ v_g = \frac{1 + w \cos^2 \theta}{\sqrt{1 + w \cos^2 \theta + \frac{m^2}{p^2}}}. \]

It is informative to evaluate (6.11) at the two extremes of \( \cos^2 \theta \); that is when \( p \) and \( \nabla \phi \) are parallel, and when they are perpendicular. If \( p \) and \( \nabla \phi \) are parallel, then the angle \( \theta \) between them vanishes and we find that the neutrino velocity is

(6.12) \[ v(\theta = 0) = \frac{1 + gu}{\sqrt{1 + \frac{1}{(1+gu)^2} \frac{m^2}{p^2}}} = \frac{c_\nu}{\sqrt{1 + \frac{1}{c_\nu^2} \frac{m^2}{p^2}}}, \]

where now \( c_\nu = 1 + gu = c_s g - g + 1 \). As in the case when \( \phi \) is uniform, if \( g = 1 \) then \( c_\nu = c_s \). We can see that the velocity of the neutrino (6.12) is almost the same as the formula given in (6.10), except that the mass term in the denominator now has a factor of \( \frac{1}{c_\nu^2} \). Here again, if the effective mass of the neutrino is zero, then the neutrino velocity is equal to the sound speed of k-essence fluctuations. This is due to the fact that the neutrino propagates on null geodesics in the emergent k-essence background. On the other hand, if \( p \) and \( \nabla \phi \) are perpendicular then

\[ v(\theta = \pi/2) = \frac{1}{\sqrt{1 + \frac{m^2}{p^2}}}, \]

and the formula for the velocity of the neutrino is the same as it would be in the absence of a k-essence field. This is because in the static field case the coupling
between the neutrino and k-essence is proportional to $p \cdot \nabla \phi$. This means that the neutrino will act as a free particle propagating on a flat Lorentzian space-time whenever it is traveling perpendicular to the direction of the field gradient. The directional dependence of the neutrino velocity in a spatially varying k-essence field stands in stark contrast to the neutrino velocity in uniform k-essence. If neutrinos traveling from a distant galaxy were to travel through a region containing a spatially varying but static k-essence field (a k-essence halo) on their way to a detector on Earth, we should expect to see evidence of anisotropy. However, even if it were possible to detect a sufficiently large neutrino flux, it is expected that any spatial variation of the k-essence field will be very small compared with its variation in time. Therefore, any anisotropy in the neutrino velocity would most likely be unobservable.

### 6.2.3 Neutrino Velocity In a General K-essence Background

Without making any assumptions about the nature of the k-essence field, the on-shell condition (6.8) becomes

$$
(1 + \frac{w\dot{\phi}^2}{2X}) \frac{E^2}{p^2} + \frac{w\dot{\phi}(\mathbf{n} \cdot \nabla \phi)}{p} \frac{E}{p} + \frac{w}{2X}(\mathbf{n} \cdot \nabla \phi)^2 - 1 = \frac{m^2}{p^2}.
$$

The solution for the $\frac{E}{p}$ is

$$
\frac{E}{p} = -\frac{w\dot{\phi}(\mathbf{n} \cdot \nabla \phi)}{2X} \pm \sqrt{1 + \frac{w}{2X}(\dot{\phi}^2 - (\mathbf{n} \cdot \nabla \phi)^2) + (1 + \frac{w\dot{\phi}^2}{2X}) \frac{m^2}{p^2}}.
$$

The choice of either a plus or minus sign in the solution reflects the two particle/anti-particle states of the neutrino: the plus sign corresponding to the neutrino, and the minus sign corresponding to the anti-neutrino. According to the Feynman-Stueckelberg interpretation of anti-particles, the anti-neutrino can be thought of as a positive energy neutrino traveling backwards in time. Thus, the solution (6.13) with the negative sign, representing the anti-neutrino energy, should have a overall
negative sign removed. Furthermore, since time is reversed, this means that in order to have the anti-neutrino traveling in the direction of $\hat{n}$, we must replace $\hat{n} \rightarrow -\hat{n}$.

In the end, the energy-momentum relation for the neutrino and anti-neutrino will be the same, and given by (6.13) with the positive sign. If we expand (6.13) to first order in $u$ and $m^2$, then $E/p$ for the (anti-)neutrino becomes

$$(6.14) \quad \frac{E}{p} \approx 1 - \frac{gu}{2X} \left( \dot{\phi} + \hat{n} \cdot \nabla \phi \right)^2 + \frac{m^2}{2p^2}. $$

From (6.13), we find that the group velocity is

$$(6.15) \quad v_g = \frac{1}{1 + \frac{w\phi^2}{2X}} \left( \frac{1 + \frac{w}{2X}(\dot{\phi}^2 - (\hat{n} \cdot \nabla \phi)^2)}{\sqrt{1 + \frac{w}{2X}(\dot{\phi}^2 - (\hat{n} \cdot \nabla \phi)^2) + (1 + \frac{w\phi^2}{2X})\frac{m^2}{p^2}}} - \frac{w}{2X} \phi (\hat{n} \cdot \nabla \phi) \right).$$

It is important to note that the neutrino velocity (6.15) does not change under a redefinition of the k-essence field variable $\phi$ unless $m^2 \neq 0$. It is easy to show that under a redefinition from $\phi$ to another field $\varphi$ defined by $\phi = g(\varphi)$, then (6.15) would remain unchanged were it not for the $m^2$ term in the denominator. Recall that $m^2$ is not the physical neutrino mass but rather a rescaled mass, which is rescaled by the conformal factor $\Omega$ in the emergent metric. After a field rescaling the effective neutrino mass becomes

$$m^2 = \frac{M^2}{\Omega^2} = \frac{M^2 L_X}{\delta} \Rightarrow m^2 = \frac{M^2}{[g'(\varphi)]^2} \frac{L_X}{\delta} = \frac{\tilde{m}^2}{[g'(\varphi)]^2}. $$

Before any objections are raised we should point out that from the beginning we have chosen a specific background k-essence field that the neutrino couples to. In essence what we have done is fix the “gauge” of the k-essence field. Therefore, it is no surprise that by changing the field variable we are changing the physics of the neutrino field. Since a field redefinition changes the conformal factor $\Omega$, a field redefinition is itself a conformal transformation. That the effective mass is the only quantity that changes under a field redefinition is a reflection of the fact that by
adding a neutrino mass we are in essence breaking the conformal invariance of the neutrino action.

6.2.4 Comparisons with Observation

In 1987 a supernova was observed [144, 145] in the Large Magellanic Cloud that provided a limited, but unique opportunity for the study of neutrino physics. During this event an increase in the background neutrino flux was detected at several neutrino observatories here on Earth. This signal was unambiguously identified as having been due to the supernova. By comparing the time interval between when the supernova was first seen and when the neutrino excess was detected, a bound on the deviation of the neutrino speed from the speed of light can be calculated. In [146] the authors found, using the available data from the supernova event, that the deviation of the neutrino speed from the speed of light can not be more than 1 part in $10^8$. In other words if $v_\nu$ is the neutrino speed then

$$\frac{|c}{v_\nu} - 1 < 10^{-8}. \quad (6.16)$$

If the neutrino is massless then $v_\nu = c_\nu$, and the left hand side of (6.16) is equal to $|gu|$ in the physically relevant static case$^2$. Thus, in the massless neutrino limit, (6.16) represents an observationally required upper bound on $|gu|$. In the most general case this bound becomes (setting $c = 1$ once again)

$$\left| \frac{1}{v_\nu} - 1 \right| < 10^{-8} \Rightarrow \frac{|gu|}{2X(\dot{\phi} + \hat{n} \cdot \nabla \phi)^2 + m^2/2E^2} < 10^{-8}. \quad (6.17)$$

The most generous upper bound that can be placed on $gu$ is $|gu| < 10^{-8}$. This can also be translated into a restriction on the k-essence sound speed $c_s$. If the k-essence field has a time-like gradient (i.e., $X > 0$), then the range of values for $c_s$ is

$$\frac{1}{1 + 10^{-8}/|g|} < c_s < \frac{1}{1 - 10^{-8}/|g|}. \quad (6.17')$$

$^2$Note that $c = 1$ in all previous sections.
As of this moment there is an insufficient amount of data to constrain $c_s$ from cosmological observables. Recent studies have shown that fits of general dark energy models to the current CMBR data are largely insensitive to the value of $c_s$ [98]. Without any other observational constraints on (or better yet, a value for) $c_s$, it is impossible to reliably estimate the value of the k-essence coupling $g$. However, with more precise data in the future it may be possible to get a better handle on the value of $c_s$. Once that has been established, (6.17) can be used to put useful restrictions on $g$.

### 6.3 Neutrino Oscillations

Experiments [147] have confirmed the phenomenon of flavor oscillations, whereby neutrinos oscillate between the possible flavor eigenstates as they travel away from their source. There are different ways of explaining this oscillation, but all mechanisms for inducing neutrino oscillations involve some term in the neutrino Lagrangian that is non-diagonal in the flavor eigenstates. Although the most popular way for inducing neutrino oscillations is by introducing a mass term [46], several other mechanisms have been proposed over the years, such as: violation of the equivalence principle (VEP) [49, 50], torsion induced neutrino oscillations [51], violation of Lorentz invariance (VLI) [52], and violation of CPT symmetry [53].

If neutrinos do indeed couple to a k-essence background in the manner we described in §6.2, then it is possible that k-essence can play a role in neutrino oscillations. There are two ways in which k-essence could affect neutrino oscillations. If the k-essence coupling is the same for each neutrino flavor, then the energy difference between energy eigenstates will not be affected. However, neutrinos coupled to the k-essence background will travel along geodesics in the emergent space-time.
This will have an affect on the phase of the neutrino wavefunction, which will be observable in the neutrino oscillation probability [149, 150, 151, 152].

Another way k-essence can influence neutrino oscillations is if the k-essence coupling $g$ is non-diagonal in the flavor eigenstate basis. Imagine now a model of two neutrino flavors that couple non-diagonally to k-essence in the flavor eigenbasis. The Lagrangian of this two neutrino system can be written as ($\alpha, \beta = 1, 2$ are flavor indices)

$$\mathcal{L} = i \sum_{\alpha=1,2} E_\alpha \bar{\nu}_K \gamma^\mu \nu_K - \frac{1}{2} \sum_{\alpha,\beta=1,2} \bar{\nu}_K (M_{\alpha\beta} E_\beta + E_\alpha M_{\alpha\beta}) \nu_K,$$

where $E_\alpha = E_{\alpha\alpha}$ and $E_{\alpha\beta} = [\text{det} E^a_{\mu}(\hat{g})]_{\alpha\beta}$. In this Lagrangian the k-essence coupling $g$ has been replaced by a matrix valued object $\hat{g}$, which is not necessarily diagonal in the flavor and mass eigenstates. We have defined $\nu_K$ as the “k-essence eigenstates”, which are the eigenstates of the k-essence coupling matrix $\hat{g}$. In general the k-essence eigenstates will not be the same as the neutrino flavor eigenstates. Because of this non-diagonal coupling of k-essence to the flavor eigenstates, the formula for the oscillation probability in the case of k-essence induced neutrino oscillations (KINO) will differ significantly from the typical mass-induced result.

In this section we will assume that the k-essence field is weakly varying, with small second derivatives so that we may effectively treat the k-essence interaction, which only involves the first derivatives of the scalar field, as a coupling constant. This greatly simplifies matters, because it allows us to diagonalize the neutrino equation of motion in the momentum space representation. This assumption is consistent with the literature where most models of k-essence take the field and its sound speed to be relatively constant in time and space. As a result of this assumption we can absorb the determinants $E_\alpha$ into a redefinition of the neutrino wavefunction $\nu_K$ and mass matrix $M_{\alpha\beta}$. Therefore, it is safe to ignore the determinant factor in our analysis.
In the first part of this section we will consider the effect of k-essence on neutrino oscillations when the k-essence coupling is equal for each flavor eigenstate.

6.3.1 Neutrino Oscillations with Flavor Diagonal K-essence Couplings

At first it may seem that when both neutrino flavors couple to k-essence identically, the usual formula for mass-induced neutrino oscillations should not change. However, even in this case the fact that neutrinos travel on geodesics in the emergent space-time implies that there will be a deviation from the flat space result. In order to best analyze neutrino oscillations in the presence of a k-essence background, we will use the simple method used in [149, 153] to study neutrino oscillations in curved space-times. In all cases, the important quantity of interest when calculating the neutrino oscillation probability is the phase of the neutrino wavefunction:

\[
|\nu_f(t, x)\rangle = e^{-i\hat{\Phi}} |\nu_f\rangle = e^{-i\int \hat{p}_\mu dx^\mu} |\nu_f\rangle.
\]

In the expression (6.18), the hats over \(\hat{\Phi}\) and the 4-momentum \(\hat{p}_\mu\) indicate that these are operators which act in the flavor space of neutrinos. If the operator \(e^{-i\hat{\Phi}}\) is non-diagonal in the flavor eigenstate basis, the result will be neutrino oscillations. The phase \(\hat{\Phi}\), written in terms of the energy and three momentum, is

\[
\hat{\Phi} = \int E\hat{I}dt - \hat{p} \cdot d\mathbf{x},
\]

where \(\hat{I}\) is the identity operator\(^3\). In a full and proper treatment of neutrino oscillations, the neutrinos must be modeled by spatially localized wave packets composed of neutrino energy eigenstates. However, if one uses this approach to study neutrino oscillations, they would find that in the relativistic limit it is acceptable to assume that both neutrino eigenstates propagate on the same null geodesic between

\(^3\text{Note, we could have assumed that the momentum operator was proportional to the identity and the energy } E \text{ was an off-diagonal operator. Both approaches are equivalent when a first order expansion in the energy is taken, as will be done in this chapter.}\)
the neutrino emitter and detector in space-time \([154, 155, 156]\)^4. Therefore the phase operator (6.19) becomes

\[
\hat{\Phi} = \int_{t_e}^{t_d} \left(E \hat{I} - \hat{p} \cdot \left[\frac{d\mathbf{x}}{dt}\right]_0\right) dt,
\]

where \(t_e\) and \(t_d\) are the values of the coordinate time at which the neutrino signal is emitted and detected, respectively. The “0” subscript on \(\frac{dx}{dt}\) denotes that this quantity is to be evaluated along a null geodesic between the emitter and detector.

To find \(\frac{dx}{dt}\), we start from the definition of the canonical momentum of a massive neutrino:

\[
p_\mu = mG_{\mu\nu} \frac{dx^\nu}{ds}.
\]

Here \(s\) is a proper time coordinate defined in the neutrino rest frame. With this, it is clear that \(\frac{dx}{dt}\) is

\[
\frac{dx}{dt} = \frac{dx}{ds} \frac{dt}{ds} = \frac{G^{i\mu}p_\mu}{G^{0\nu}p_\nu} = \frac{\mathbf{p} - \frac{w}{2X}(\dot{\phi}E + \mathbf{p} \cdot \nabla \phi)\nabla \phi}{E + \frac{w}{2X}(\dot{\phi}E + \mathbf{p} \cdot \nabla \phi)}.
\]

Along a null geodesic the energy and momentum of a massless neutrino are denoted by \(E^{(0)}\) and \(p^{(0)}\), respectively. Let us denote the unit vector in the direction of the neutrino momentum by \(\hat{n}\). Thus, \(\hat{n} \cdot \left[\frac{dx}{dt}\right]_0\) is given by

\[
\hat{n} \cdot \left[\frac{dx}{dt}\right]_0 = -\frac{w}{2X} \dot{\phi}(\hat{n} \cdot \nabla \phi) + \sqrt{1 + \frac{w}{2X}(\dot{\phi}^2 - (\hat{n} \cdot \nabla \phi)^2)} = \frac{E^{(0)}}{p^{(0)}}.
\]

To relate \(E\) and \(\mathbf{p}\), recall the discussion from section §6.2.3. It is easy to show from the work done there that the momentum as a function of \(E\) expressed to leading

\[^4\text{In the wave packet treatment, an exponential damping is found in the final result for the oscillation probability in the relativistic limit. This is due to the decoherence of the superposition of neutrino wavefunctions. So far there has been no solid evidence for decoherence effects [157], and we are therefore safe in ignoring this possibility in our analysis.}\]
order in $m^2$ is

$$p \approx E \frac{\frac{w}{2E} \dot{\phi} (\hat{n} \cdot \nabla \phi) + \sqrt{1 + \frac{w}{2E} (\dot{\phi}^2 - (\hat{n} \cdot \nabla \phi)^2)}}{1 - \frac{w}{2E} (\hat{n} \cdot \nabla \phi)^2} - \frac{m^2}{2E} \sqrt{1 + \frac{w}{2E} (\dot{\phi}^2 - (\hat{n} \cdot \nabla \phi)^2)}.$$  

Now remember that $m^2 = \Omega^{-2} M^2$. In order for there to be neutrino oscillations the flavor and mass eigenstates must not be the same. Therefore, we have to replace the $M^2$ with a matrix $\hat{M}^2$ that is non-diagonal in the flavor eigenstate basis. Since only phase differences between flavor eigenstates will be important, we can ignore terms not proportional to $\hat{M}^2$ in the phase operator (6.20). With (6.21) and (6.22) the expression for the phase operator (6.20), modulo terms proportional to the identity, becomes

$$\hat{\Phi} = \frac{\hat{M}^2 L(\ell)}{2E^{(0)}},$$  

where $x_e$ and $x_d$ are the neutrino and detector positions, respectively. We have defined an effective distance $L$, which is a function of the coordinate distance $\ell = |x_d - x_e|$ between the neutrino emitter and detector. The expression for $L(\ell)$ is given by

$$L(\ell) = \int_{x_e}^{x_d} \frac{\Omega^{-2} (\hat{n} \cdot dx)}{\sqrt{1 + \frac{w}{2E} (\dot{\phi}^2 - (\hat{n} \cdot \nabla \phi)^2)}}.$$  

In general this distance will differ from the true emitter/detector separation due to the neutrino coupling to the k-essence medium. There are two sets of neutrino eigenstates in this system: the flavor eigenstates $|\nu_f\rangle = \{|\nu_e\rangle, |\nu_\mu\rangle\}$, which couple diagonally to the weak current; and the mass eigenstates $|\nu_m\rangle = \{|\nu_{m_1}\rangle, |\nu_{m_2}\rangle\}$ that define the basis in which the mass matrix is diagonal. In general these two sets will not be equivalent. However, we can define an $SU(2)$ matrix $V$ such that

$$|\nu_f\rangle = \sum_m V_{fm} |\nu_m\rangle.$$
The matrix $V$ is referred to as the mass mixing matrix. The most general $SU(2)$ matrix can be represented in terms of an angle and three phases:

$$e^{i\chi} \begin{bmatrix} e^{-i\alpha} & 0 \\ 0 & e^{i\alpha} \end{bmatrix} \begin{bmatrix} \cos \theta & \sin \theta \\ -\sin \theta & \cos \theta \end{bmatrix} \begin{bmatrix} e^{-i\beta} & 0 \\ 0 & e^{i\beta} \end{bmatrix}.$$

In general, $V$ is an $SU(2)$ matrix and will therefore have three additional phase degrees of freedom. However, in this case we can absorb these phases into a redefinition of the neutrino wavefunction. Thus, we can safely ignore the phases in $V$, and as a result the mixing between the flavor and mass eigenstates is determined by a single angle, which we will call $\theta_M$:

$$V = \begin{bmatrix} \cos \theta_M & \sin \theta_M \\ -\sin \theta_M & \cos \theta_M \end{bmatrix}.$$

Written in the flavor eigenstate basis, the phase operator (6.23) is

$$\hat{\Phi}(\ell) = \frac{L(\ell)}{2E} V \begin{bmatrix} m_1^2 & 0 \\ 0 & m_2^2 \end{bmatrix} V^\dagger = \frac{\Delta m^2 L(\ell)}{4E} \begin{bmatrix} -\cos 2\theta_M & \sin 2\theta_M \\ \sin 2\theta_M & \cos 2\theta_M \end{bmatrix} + \frac{\bar{m}^2 L(\ell)}{2E} \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix},$$

where $m_1$ and $m_2$ are the masses of the two neutrino mass eigenstates, and $\Delta m^2 = m_2^2 - m_1^2$ and $\bar{m}^2 = \frac{m_3^2 + m_1^2}{2}$. Note that in the interest of simplicity, we have dropped the “(0)” superscript on the energy $E$. Since we are interested in finding probabilities, we can subtract from the phase matrix any term proportional to the identity matrix without changing our final result. Therefore, we will ignore the very last term in (6.24) from here on out, and thus, the phase factor in (6.18) becomes

$$e^{-i\phi} = -i \sin \left( \frac{\Delta m^2 L(\ell)}{4E} \right) \begin{bmatrix} -\cos 2\theta_M & \sin 2\theta_M \\ \sin 2\theta_M & \cos 2\theta_M \end{bmatrix} + \cos \left( \frac{\Delta m^2 L(\ell)}{4E} \right) \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix}.$$
The factor $e^{-i\hat{\Phi}}$ plays the role of a evolution operator. Denote by $|\nu_f\rangle$ the initial state of a neutrino of flavor $f$. The time at which the neutrino is created can be defined, without loss of generality, as occurring at $t = 0$. The wavefunction of this neutrino at some later time $t > 0$, after it has traveled a distance $\ell$ away from its source, will be given by the action of $e^{-i\hat{\Phi}(\ell)}$ on the neutrino wavefunction $|\nu_f\rangle$:

$$|\nu_f(\ell)\rangle = \sum_{f'} [e^{-i\hat{\Phi}(\ell)}]_{ff'} |\nu_{f'}\rangle.$$ 

Applying the phase operator (6.25) to the neutrino ket $|\nu_f\rangle$, we find that the flavor eigenstate wavefunctions at a later time $t > 0$, when the neutrinos have traveled a distance $\ell$, are

$$|\nu_e(\ell)\rangle = i\sin\left(\frac{\Delta m^2 L(\ell)}{4E}\right) (\cos 2\theta_M |\nu_e\rangle - \sin 2\theta_M |\nu_\mu\rangle) + \cos\left(\frac{\Delta m^2 L(\ell)}{4E}\right) |\nu_e\rangle,$$

$$|\nu_\mu(\ell)\rangle = -i\sin\left(\frac{\Delta m^2 L(\ell)}{4E}\right) (\sin 2\theta_M |\nu_e\rangle + \cos 2\theta_M |\nu_\mu\rangle) + \cos\left(\frac{\Delta m^2 L(\ell)}{4E}\right) |\nu_\mu\rangle.$$

Suppose a neutrino of flavor $f$ is created from some source of interest (e.g., the sun, atmosphere, nuclear reactor, etc.) at some time $t_e$. The probability of this neutrino appearing as an $f'$ flavored neutrino in a detector here on Earth, at some later time $t_d$, after having traversed a distance $\ell$, is given by the expression

$$(6.26) \quad P(\nu_f \rightarrow \nu_{f'}) = |\langle \nu_{f'} | \nu_f(\ell) \rangle|^2 = \delta_{f, f'} + (-1)^{\delta_{f, f'}} \sin^2 2\theta_M \sin^2\left(\frac{\Delta m^2 L(\ell)}{4E}\right).$$

As we can see here the formula in (6.26) is almost exactly the same result one gets for the oscillation probability for mass-induced neutrino oscillations in flat space. The only difference is that the effective distance $L$ takes the place of the coordinate distance $\ell$.

6.3.2 Neutrino Oscillations with Flavor Non-diagonal K-essence Couplings

At the beginning of this chapter, we introduced a neutrino/k-essence coupling parameter that we denoted by $g$. In the last subsection we assumed that this coupling
was the same for each neutrino flavor. In this subsection, however, we will consider what happens when this coupling $g$ becomes a matrix valued object $\hat{g}$ that operates in the neutrino flavor space. In this case the equation for the energy-momentum relation (6.14) becomes a matrix equation whose eigenvalues represent the energies of the neutrino energy eigenstates, which in general are not the same as the flavor, mass, and k-essence eigenstates. A formula for the momentum as a function of energy can be derived from the information found in section §6.2.3. If we expand this formula around $g = 0$ and $m^2 = 0$ to leading order, and then replace $g$ and $M^2 = \Omega^2 m^2$ with matrices that act in the neutrino flavor space, we end up with

$$\hat{p} \approx E \hat{I} + \frac{E \hat{G}}{2} - \frac{\hat{M}^2 \Omega^{-2}}{2E},$$

where $\hat{G} = \frac{\dot{\phi}}{\mathcal{N}} (\dot{\phi} + \hat{n} \cdot \nabla \phi)^2$. Here $\hat{G}$ and $\hat{M}^2$ are operators in the flavor space. In general it will not be possible to diagonalize $\hat{G}$ and $\hat{M}^2$ at the same time, since there is no reason to assume that the k-essence and mass eigenstates are the same. We can relate the different sets of eigenstates with two $SU(2)$ matrices. Define two $SU(2)$ matrices $U$ and $V$ such that

$$|\nu_f\rangle = \sum_m V_{fm} |\nu_m\rangle, \quad |\nu_f\rangle = \sum_\alpha U_{f\alpha} |\nu_\alpha\rangle \Rightarrow |\nu_\alpha\rangle = \sum_{f,m} U_{f\alpha}^\dagger V_{fm} |\nu_m\rangle.$$

In the case where there are three sets of neutrino eigenstates, one can not simply disregard the phases in the mixing matrices $U$ and $V$ as can be done in the standard treatment of purely mass-induced neutrino oscillations. While we can eliminate most of the phases through a redefinition of the neutrino eigenstates, there is not enough freedom to get rid of them all. After any redefinition of the neutrino eigenstates, there will still be a single overall phase left. Let’s call this residual phase $\alpha$ and
define the k-essence and mass mixing matrices as

\[
U = \begin{pmatrix}
e^{-i\alpha} \cos \theta_K & e^{-i\alpha} \sin \theta_K \\
-e^{i\alpha} \sin \theta_K & e^{i\alpha} \cos \theta_K
\end{pmatrix}, \quad V = \begin{pmatrix}
\cos \theta_M & \sin \theta_M \\
-\sin \theta_M & \cos \theta_M
\end{pmatrix}.
\]

If \(G_1\) and \(G_2\) are the eigenvalues of \(\hat{G}\), the momentum operator can be written as

\[
\hat{p} \approx E \begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix} + \frac{E}{2} U \begin{pmatrix} G_1 & 0 \\ 0 & G_2 \end{pmatrix} U^\dagger - \frac{1}{2E} V \begin{pmatrix} \tilde{m}_1^2 & 0 \\ 0 & \tilde{m}_2^2 \end{pmatrix} V^\dagger,
\]

where \(\tilde{m}_i^2 = \Omega^{-2} m_i^2\). If we substitute this for the momentum operator in the integrand in (6.20) we find that, modulo terms proportional to the identity matrix, we get

\[
E\hat{I} - \hat{p} \cdot \left[ \frac{dx}{dt} \right]_0 = \frac{\Delta\tilde{m}^2}{4E} \begin{pmatrix} -\cos 2\theta_M & \sin 2\theta_M \\ \sin 2\theta_M & \cos 2\theta_M \end{pmatrix} - \frac{E\Delta G}{4} \begin{pmatrix} -\cos 2\theta_K & e^{-2i\alpha} \sin 2\theta_K \\ e^{2i\alpha} \sin 2\theta_K & \cos 2\theta_K \end{pmatrix}
\]

(6.27)

where \(\Delta\tilde{m}^2 = \tilde{m}_2^2 - \tilde{m}_1^2\). Here we have introduced a new parameter \(y\), which is defined as \(y = -\frac{E^2 \Delta G}{\Delta m^2}\), where \(\Delta G = G_2 - G_1\). This operator can be written in a much more compact form by defining new variables:

\[
\frac{\Delta M^2}{4E} \cos 2\theta_L = \frac{\Delta\tilde{m}^2}{4E} (\cos 2\theta_M + y \cos 2\theta_K),
\]

\[
\frac{\Delta M^2}{4E} e^{-i\sigma} \sin 2\theta_L = \frac{\Delta\tilde{m}^2}{4E} (\sin 2\theta_M + y \sin 2\theta_K e^{-2i\alpha}).
\]

Note that the phase \(\sigma\), while present here, will not affect our final result. Inverting these:

\[
\sin^2 2\theta_L = \frac{\sin^2 2\theta_M + 2y \sin 2\theta_M \sin 2\theta_K \cos 2\alpha + y^2 \sin^2 2\theta_K}{1 + 2y \cos 2\Theta + y^2},
\]

\[
\tan \sigma = \frac{\sin 2\theta_K \sin 2\alpha}{\sin 2\theta_M + y \sin 2\theta_K \cos 2\alpha};
\]

\[
\frac{\Delta M^2}{4E} = \frac{|\Delta\tilde{m}^2|}{4E} \sqrt{1 + 2y \cos 2\Theta + y^2},
\]

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where

\[
\cos 2\Theta = \cos 2\theta_M \cos 2\theta_K + \sin 2\theta_M \sin 2\theta_K \cos 2\alpha.
\]

With these, equation (6.27) can be written as

\[
(6.28) \quad E\hat{I} - \hat{p} \cdot \left[ \frac{dx}{dt} \right]_0 = \frac{\Delta M^2}{4E} \begin{bmatrix}
-\cos 2\theta_L & e^{-i\sigma} \sin 2\theta_L \\
e^{i\sigma} \sin 2\theta_L & \cos 2\theta_L
\end{bmatrix}.
\]

The phase operator is therefore

\[
(6.29) \quad e^{-i\Phi(\ell)} = -i \sin \varphi(\ell) \begin{bmatrix}
-\cos 2\theta_L & e^{-i\sigma} \sin 2\theta_L \\
e^{i\sigma} \sin 2\theta_L & \cos 2\theta_L
\end{bmatrix} + \cos \varphi(\ell) \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix},
\]

where \(\varphi(\ell) = \int_{x_e}^{x_d} \frac{\Delta M^2}{4E} \, dx\) and \(\ell = x_d - x_e\). Using (6.29) as our neutrino evolution operator, the flavor eigenstates at a later time after the neutrino has traveled a distance \(\ell\), are given by:

\[
|\nu_e(\ell)\rangle = i \sin \varphi(\ell) \left( \cos 2\theta_L |\nu_e\rangle - e^{-i\sigma} \sin 2\theta_L |\nu_\mu\rangle \right) + \cos \varphi(\ell) |\nu_e\rangle,
\]

\[
|\nu_\mu(\ell)\rangle = -i \sin \varphi(\ell) \left( e^{i\sigma} \sin 2\theta_L |\nu_e\rangle + \cos 2\theta_L |\nu_\mu\rangle \right) + \cos \varphi(\ell) |\nu_\mu\rangle.
\]

Therefore, if a neutrino of flavor \(f\) is created and travels a distance \(\ell\) to a detector, the probability that the neutrino is observed as an \(f'\) flavored neutrino is

\[
P(\nu_f \rightarrow \nu_{f'}) = |\langle \nu_{f'} | \nu_f(\ell) \rangle|^2 = \delta_{f,f'} + (-1)^{\delta_{f,f'}} \sin^2 2\theta_L \sin^2 \varphi(\ell).
\]

It is interesting to consider the case when the k-essence field does not vary rapidly in space and time compared to terrestrial scales. In that case we can treat the integrand as a constant, and therefore the phase becomes

\[
\varphi(\ell) = \frac{\Delta \tilde{m}^2}{4E} \ell \sqrt{1 + 2y \cos 2\Theta + y^2}.
\]
Since the phase is proportional to the distance between the neutrino emitter and detector, we can write the oscillation probability as

\[ P(\nu_f \rightarrow \nu_{f'}) = \delta_{f,f'} + (-1)^{\delta_{f',f}} \sin^2 2\theta_L \sin^2 \frac{\pi l}{\lambda}, \]

where \( \lambda \) is the oscillation length, which is defined as

\[ \lambda = \frac{\pi \ell}{\varphi(\ell)} = \frac{4\pi E}{|\Delta \tilde{m}^2|} \frac{1}{\sqrt{1 + 2y \cos 2\Theta + y^2}}. \]

It is interesting to compare neutrino oscillations that are either induced purely by k-essence or by mass. In the case where neutrino oscillations are due to mass entirely, the inverse oscillation length goes like

\[ 4\pi \lambda^{-1} = \frac{|\Delta \tilde{m}^2|}{E}. \]

Likewise, if the mass of neutrinos vanishes, then the inverse oscillation length goes like

\[ 4\pi \lambda^{-1} = E |\Delta G| = E \left| \frac{(g_2 - g_1)u}{X} \right| (\dot{\phi}^2 + \hat{n} \cdot \nabla \phi)^2. \]

Comparing (6.31) and (6.32) it is apparent that neutrino oscillations induced by either a flavor non-diagonal mass term, or a flavor non-diagonal k-essence coupling will lead to noticeably different energy dependences for the oscillation length. If neutrino oscillations are entirely induced by k-essence then the oscillation length goes like \( \lambda^{-1} \sim E \). The result (6.32) should be compared to neutrino oscillations induced by either the VLI or VEP mechanisms. Both the VLI and VEP mechanisms have the same \( \lambda^{-1} \sim E \) behavior that the KINO mechanism has. This should come as no surprise since the flavor-dependent emergent metric \( G_{\mu\nu}^{(\alpha)} \) can be viewed as a regular space-time metric but with a flavor-dependent, and therefore equivalence principle violating, gravitational constant. Equations (6.30) and (6.32) immediately
tell us that in order for there to be flavor oscillations not only must $g_1 \neq g_2$, but also $u \neq 0$. Since $u$ vanishes for $c_s = 1$, this means that there must be Lorentz violation if k-essence is to have any effect on neutrino oscillations. K-essence, therefore, acts as a Lorentz violating aether background, and can be the motivation behind models of neutrino oscillation that invoke Lorentz violation.

Analysis of the available data from current and past neutrino observatories have tended to favor mass-induced neutrino oscillations, which can produce the desired $\lambda \propto E$ type behavior. Therefore, the KINO mechanism alone does not suffice to explain the observations seen in the numerous neutrino experiments that have been carried out. However, it is still possible that k-essence could be a subleading contribution to neutrino oscillations, with mass being the dominant cause. Studies have looked into the possibility of alterations to the leading order $\lambda \propto E$ dependence, and have been able to place very tight constraints on the coefficients of subleading contributions to the energy dependence of $\lambda$. In [158] they considered the possibility of different mechanisms inducing oscillations in the $\nu_\mu \leftrightarrow \nu_\tau$ channel, among which were the VLI and VEP mechanisms. Since the VLI and VEP scenarios both lead to the same energy dependence for the oscillation length that KINO does, the constraints on the VLI and VEP coefficients can be easily translated into a bound on $\Delta G$:

$$|\Delta G| < 1.2 \times 10^{-23}.$$ 

This is the most conservative bound that can be placed on $|\Delta G|$, and is independent of the mixing angle $\theta_K$. This bound would seem to cast doubt on KINO as an even subleading effect in the $\nu_\mu \leftrightarrow \nu_\tau$ channel. We would like to emphasize that this is a strong indication of the nature of the coupling of neutrinos to the k-essence scalar. Our analysis suggests that if the k-essence scalar field exists, in order to be
phenomenologically viable, its couplings must be flavor blind.

Although KINO would seem to be immediately discounted from consideration, it may be possible to realistically consider this model if we are ready to include further symmetry violating terms in our action. It has been shown in [159] that if both Lorentz, and certain types of CPT violating terms are included in the neutrino action, then it becomes possible to create pseudo-mass terms at high energies just when Lorentz violating effects should be taking over. These types of models open up the possibility of a unified explanation for all the existing neutrino data, including the controversial LSND results [160]. Since these models have been at least qualitatively compatible with experiment, this leaves open the possibility that k-essence could still play a role in neutrino oscillations.

As we have just mentioned, this approach requires a specific kind of CPT violating term in the neutrino action. Although it is always possible to put terms into the action arbitrarily, in principle it might be possible for k-essence to be a source for these as well. K-essence could lead to the needed CPT breaking terms in the neutrino action by one of two ways. The first is by a possible axial vector term in the spinor covariant derivative. In this chapter we have assumed that space-time is flat, and as a result the covariant derivative is proportional to $\gamma^\mu$. However, if more general curved space-times are considered then in general the spinor connection will have a nonzero axial vector part. Another way in which k-essence can source the necessary CPT violating term is by considering the possibility of torsion in the emergent space-time. Such a treatment of k-essence would require taking into account the spin-orbit coupling between the neutrino and the k-essence field, and treating the metric and connection of the emergent space-time as independent variables each requiring their

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5See appendix §B.
own field equations. Torsion in k-essence backgrounds is an interesting possibility and deserves further study in its own right.

6.4 Conclusion

In this chapter we have investigated the effects on neutrino velocity and oscillations when neutrinos couple to the emergent metric $G_{\mu \nu}$ created by a background k-essence field. Specifically, we have studied the results of replacing the vierbein $e_\alpha^\mu$ of the gravitational background with the vierbein $E_\alpha^\mu$ of the emergent k-essence geometry in the neutrino action. The first implication of this coupling is a change in the neutrino dispersion relation, which means that the neutrino velocity is dependent on the k-essence background in which it propagates. Without k-essence, massless neutrinos always propagate at the speed of light with respect to local observers. However, with k-essence, massless neutrinos will move at a new speed $c_\nu$ with respect to an observer who perceives a uniform k-essence background. Therefore, if future observations show that neutrinos do travel at less than the speed of light, it can not be determined conclusively if this is due to neutrinos being massive, or if neutrinos are massless but coupled to k-essence. However, at present, no observations have found a measurable difference between the speed of light and neutrinos. If neutrinos are assumed to be massless, then past observations constrain the deviation of $c_\nu$ from the speed of light to no more than 1 part in $10^8$.

The other effect that k-essence has on neutrinos is in the phenomenon of neutrino oscillations. We have shown that if the neutrino/k-essence coupling $g$ is promoted to the status of an operator $\hat{g}$ that acts in the neutrino flavor space, neutrino oscillations are produced even in the absence of a neutrino mass term. We found that if neutrino oscillations are caused entirely by k-essence, then the oscillation length depends on
the neutrino energy like $\lambda \sim E^{-1}$. This is to be contrasted with the result from mass-induced neutrino oscillations where $\lambda \sim E$. Thus, while neutrino oscillations can be induced by k-essence, it will lead to a drastically different energy dependence for the oscillation length. Current data seems to favor a $\lambda \sim E$ behavior, which implies an important constraint on the couplings of neutrinos to the k-essence scalar. In order to be phenomenologically viable these couplings must be flavor diagonal. Our analysis is a very good example of how present observations can be used to constrain the form and magnitude of couplings between dark energy and visible matter.

Neutrino oscillations induced by k-essence have many of the same properties that some [49, 50, 52] earlier proposed mechanisms had, and in fact k-essence can be seen as a realization of these past phenomenologically motivated models. In VLI models, the Lorentz violation is incorporated by introducing some preferred 4-vector into the neutrino action. This preferred 4-vector has the interpretation in k-essence as the 4-gradient of the k-essence field. In fact, as can be seen from the formulae in section §6.3, the physical effects vanish at $c_s = 1$ where our model has exact Lorentz invariance. The VEP mechanism attempts to explain the source of neutrino oscillations as a consequence of a flavor non-diagonal coupling to gravity. Although perhaps correct from a theoretical point of view, this mechanism calls for a very drastic change in our understanding of fundamental physics; namely it requires us to give up the long held notion of the equivalence principle. The type of coupling studied in this chapter mimics exactly the VEP mechanism, but since the flavor non-diagonal coupling is in the k-essence sector instead of the gravitational sector, the equivalence principle is maintained. In short, k-essence can be seen as the source of Lorentz violation in the VLI model, and alternatively, a reinterpretation of the VEP model. KINO is therefore a theoretically and phenomenologically motivated manifestation
of earlier attempts at alternatives to mass-induced neutrino oscillations.

In this chapter, we have coupled neutrinos to the k-essence scalar using the emergent metric which is covariantly constant. In a future work we will show how this can be extended to more general couplings by introducing the analogue of torsion for the emergent space-time. This generalization will allow us to analyze the role, if any, of CPT violation in the coupling of k-essence to fermions. If CPT violation can be naturally produced, then it would be possible to discuss more realistic models of neutrino oscillations as discussed in [159].
CHAPTER VII

Conclusion

The work in this thesis has primarily concerned itself with non-canonical scalar fields and their uses in cosmology. As we have seen, such theories have a number of potential applications in cosmology, many of these dealing with some of the most important questions we have about our universe. In chapters §II and §III we studied the possibility that k-essence is the source of dark matter. In these chapters we found general conditions that must be satisfied in order for k-essence halos to be stable. In chapter §II, we looked at k-essence halo solutions in flat space-time. The main point of this chapter was that stable k-essence halos are possible, however, it requires that the energy density must be negative in a small region around the center of the halo. Despite this, the total energy of the halo is positive, which means that the gravitational interactions between galaxies with k-essence halos are consistent with observations. In chapter §III we extended this analysis to account for the gravitational back-reaction that a large k-essence halo would create. In this chapter we studied the consequences of k-essence halos both with and without a black hole at the center of the condensate. It was found that with no central black hole, the conditions for stable k-essence halos found in chapter §II also applied when back-reaction effects are considered. When a black hole was placed in the center it was
found that stable solutions were possible if there is a shell of negative energy density at the horizon of the black hole, or if the space-time is not asymptotically flat. Again it was shown that although there was a region of negative energy density, the total energy of the halo was positive both when there was a black hole and when there wasn’t. If one considers negative energy densities unacceptable, then our results can be considered a no-go theorem for k-essence halos. However, results from string theory suggest that negative energy densities are not necessarily unphysical [71]. Therefore, our analysis suggests that k-essence halos are at least possible from the standpoint of classical physics. However, to say so definitively it must be shown that such solutions are stable against quantum mechanical vacuum decay. This question will have to be the subject of future work.

Chapter §IV is an outgrowth of the work done in chapter §III. The possibility of k-essence condensates around black holes raises the question of whether such configurations are stable due to the tendency of black holes to “eat” any material around them. In studying this question, one usually considers solutions that are stationary; i.e., solutions where the energy density, pressure, and four velocity are constant in time. Thus, it is natural to ask, under what conditions are stationary solutions possible in k-essence? Chapter §IV was concerned with finding these conditions, and found that stationary solutions are not the rule, but rather the exception. In that chapter we showed that stationary k-essence solutions around black holes are only possible when the action of the scalar field is equivalent to the kinetic k-essence action after a field redefinition. This implies that all actions that admit stationary solutions are invariant under constant field shifts: $\phi \rightarrow \phi + \text{const}$.

In chapter §V we laid out a method for using data on the scalar and tensor perturbation spectra and the non-gaussianity to find the inflationary action, even
if this action is non-canonical. The problem involved finding the function $L(X, \phi)$, given the three functions $P_s(k)$, $P_h(k)$ and $f_{NL}(k)$, which represent the scalar, tensor, and non-gaussianity spectra, respectively. However, the data are single variable functions of the scale $k$; these alone are not enough to determine $L(X, \phi)$, which is a function of two variables. Therefore, we assumed that a user of the algorithm specifies a “partition” of the action. We defined the partition of an action $L(X, \phi)$ as a multivariable function $P$ such that

$$L(X, \phi) = P(g_1(X), ..., g_m(X), f_1(\phi), ..., f_n(\phi)).$$

The form of the function $P$ is assumed by the user, while the single variable functions $f_\alpha$ and $g_\beta$ are unknown. The procedure we laid out in chapter §V gives the user a system of equations that can be used to find (at most) three of the $f$’s and $g$’s once we are given the functional forms of $P_s(k)$, $P_h(k)$ and $f_{NL}(k)$. If there are less than three unknown functions, the reconstruction equations that are not used establish consistency relations between the CMBR observables. We applied this procedure ourselves to various different partitions $P$. The example we concentrated the most on was a generalized version of the DBI inflationary action. This was

$$L_{DBI}(X, \phi) = -f^{-1}(\phi)\sqrt{1 - 2f(\phi)X} + f^{-1}(\phi) - V(\phi).$$

The functions $f$ and $V$ are assumed to be unknown and the objective was to find them using the data. Since there were only two unknown functions only the scalar and tensor data were needed to find $f$ and $V$. Using our procedure we found that physically acceptable inflationary actions did not exist for all functions of $P_s$ and $P_h$. What’s more, we showed that if $P_s$ is nearly scale-invariant, $f$ and $V$ are approximately what theoretical models suggest they should be.

In the final chapter we considered the effects that a k-essence field might have
on neutrinos. In chapter §VI we proposed that neutrinos couple to the k-essence induced metric $G_{\mu\nu}$, rather than the space-time metric $g_{\mu\nu}$. One effect that this had on the neutrinos was it modified the energy-momentum relation. In the presence of this coupling, the neutrino velocity was in general different from the speed of light, even without mass. In a frame where the k-essence is spatially uniform, the motion of neutrinos is the same as any free particle, but their speed is now related to the k-essence sound speed. Later in the same chapter we showed that k-essence could also induce neutrino oscillations even without a neutrino mass term. It was shown that neutrino oscillations induced purely by k-essence led to an oscillation length that went like $\lambda \sim E^{-1}$. This conflicts with the case of purely mass induced neutrino oscillations, which result in a $\lambda \sim E$ type behavior. Thus, k-essence induced neutrino oscillations have a very different observational signature than neutrino oscillations created by mass. However, observations favor a leading order $\lambda \sim E$ behavior. Therefore, our results put tight constraints on the magnitude and form of neutrino/k-essence interactions.

The results discussed in this thesis show that non-canonical scalar fields have a number of different applications in cosmology. Although our work has helped to elucidate these potential uses, a number of avenues for further research remain. In the area of k-essence dark matter, questions remain about the stability of k-essence halos under quantum fluctuations. Since supermassive black holes are thought to be at the center of every galaxy, it is important to understand how k-essence halos evolve in the presence of black holes. As we found in chapter §IV, stationary solutions are in general not possible. Studying this problem will no doubt require sophisticated numerical simulations.

Since k-essence can be constructed to behave as either dark matter or dark energy,
it is tempting to ask: “is there a k-essence theory that behaves like dark matter at small scales and dark energy at large scales?” If such theories exist it would be an exciting new way of unifying dark matter and dark energy into a single framework. Furthermore, if k-essence and k-inflation can be unified, then we could have a single theory that can explain inflation, dark matter and dark energy. Finally, if k-essence does exist it is important to understand the possible interactions it can have with matter. Such interactions could be used to detect and understand k-essence in particle accelerators here on Earth. Wherever the research leads us next, k-essence has certainly opened up new possibilities, and it warrants our further attention.
APPENDICES
In this appendix we will discuss the issues of causality and the Cauchy problem in k-essence theories. In k-essence theories, one is faced with the problem of solving the equations of motion of a field that lives on some background geometry defined by the k-essence induced metric $G_{\mu\nu}$. As we will see in this appendix, certain conditions on the k-essence field and its action must be satisfied in order for this problem to be well posed.

To start, let $\phi$ be a k-essence field with a Lagrangian of the general form $L = L(X, \phi)$, where $X = \frac{1}{2} g_{\mu\nu} \nabla_\mu \phi \nabla_\nu \phi$. The equation of motion for a k-essence field in a general gravitational background $g_{\mu\nu}$ is given by
\begin{equation}
\tilde{G}^{\mu\nu}[\phi] \nabla_\mu \phi \nabla_\nu \phi = L_\phi - 2X L_X \phi, \quad \text{where} \quad \tilde{G}^{\mu\nu}[\phi] = L_X g^{\mu\nu} + L_{XX} \nabla^\mu \phi \nabla^\nu \phi.
\end{equation}
(A.1)

where $L_X = \frac{\partial L}{\partial X}$, $L_\phi = \frac{\partial L}{\partial \phi}$, etc. In (A.1), $\tilde{G}^{\mu\nu}$ plays the role of an effective metric in which the k-essence field propagates. The energy-momentum tensor $T_{\mu\nu}$ of a general k-essence field is

\begin{equation}
T_{\mu\nu} = L_X \nabla_\mu \phi \nabla_\nu \phi - g_{\mu\nu} L.
\end{equation}

From the energy-momentum tensor, we can define the energy density and pressure of the k-essence field depending on whether $\nabla_\mu \phi$ is time-like ($X > 0$) or space-like
\( X < 0 \):

\[
\rho = 2XL_X - L, \quad p = L; \quad \text{time-like}
\]

\[
\rho = -L, \quad p = L - 2XL_X; \quad \text{space-like}
\]

In either case, the null energy condition (i.e., \( \rho + p > 0 \)), which we assume throughout this thesis, demands that \( L_X > 0 \).

The behavior of small perturbations \( \pi \) in the k-essence field \( \phi \) around some background \( \phi_0 \) are described by the equation of motion for a canonical scalar field:

\[
G^{\mu\nu}[\phi_0]D_\mu D_\nu \pi + M_{eff}^2[\phi_0] \pi = 0,
\]

where

\[
G^{\mu\nu}[\phi_0] = \frac{\delta}{L_X^2} \tilde{G}^{\mu\nu}[\phi_0], \quad M_{eff}^2 = \frac{\delta}{L_X^2} \left( 2XL_X \phi_0 - L_{\phi\phi} + \frac{\partial \tilde{G}^{\mu\nu}}{\partial \phi} \nabla_\mu \nabla_\nu \phi_0 \right).
\]

Here, \( \delta \) is a new parameter defined by \( \delta^2 = \frac{L_X}{L_X + 2XL_{XX}} \). The parameter \( \delta \) is equal to the sound speed of k-essence fluctuations in the case where \( \nabla_\mu \phi_0 \) is time-like. By definition, the sound speed \( c_s \) is given by

\[
c_s^2 = \frac{\partial p/\partial X}{\partial \rho/\partial X} \bigg|_\phi = \begin{cases} 
\frac{L_X}{L_X + 2XL_{XX}}; & X > 0 \\
\frac{L_X + 2XL_{XX}}{L_X}; & X < 0
\end{cases}
\]

In either case, absence of superluminal propagation demands that \( L_{XX} \geq 0 \). In this thesis we have limited our study mostly to subluminal k-essence theories. However, it has been shown that even with superluminal sound speeds, causality can still be maintained [31].

The question of whether or not one can solve the equations of motion for \( \phi_0 \) and \( \pi \) is not trivial, but depends on the k-essence action and the initial conditions of the field. The problem of finding solutions to a system of equations given the initial conditions is known as the initial value problem, or Cauchy problem. The existence of well behaved solutions hinges on whether or not \( G^{\mu\nu} \) has a well defined casual
structure, with no closed time-like curves. In order to say more, let us define a space-time \((\mathcal{M}, g_{\mu\nu})\) as a set consisting of a differentiable manifold \(\mathcal{M}\), equipped with a metric tensor \(g_{\mu\nu}\), whose signature is \((+,-,-,-)^1\). It is known [161] that a space-time is \textit{causally stable}; i.e., absent of closed time-like curves, if and only if there exists some function \(f\) on \(\mathcal{M}\) such that \(\nabla^\mu f\) is a future directed, time-like vector field. The function \(f\) can be thought of as the time coordinate.

In order to say whether or not the metric \(G_{\mu\nu}\) is causally stable, \(G_{\mu\nu}\) has to be \textit{hyperbolic}; i.e., it must have a Lorentzian signature. Unless \(G_{\mu\nu}\) is hyperbolic, one can not define the notion of time-like and space-like vectors. Assuming that \(g_{\mu\nu}\) is Lorentzian, this amounts to requiring that \(G_{\mu\nu}\) be such that

\[
\det G_{\mu\nu} < 0 \implies 1 + 2X \frac{L_{XX}}{L_X} > 0.
\]

Note that the hyperbolicity condition is equivalent to the condition for the stability of k-essence perturbations: \(c_s^2 > 0\). If \(G_{\mu\nu}\) satisfies this condition, let’s assume that \((\mathcal{M}, g_{\mu\nu})\) is a causally stable space-time. Therefore, we can define a function \(t\) that represents the time in the appropriate coordinate system, and is such that \(n_\mu = \partial_\mu t\) is a future directed, time-like (with respect to \(g^{\mu\nu}\)) vector field. In order for this vector to be time-like with respect to \(G^{\mu\nu}\), then

\[
(A.2) \quad G^{\mu\nu} n_\mu n_\nu = \frac{\delta}{L_X} \left(n^2 + \frac{L_{XX}}{L_X} (n_\mu \nabla^\mu \phi_0)^2\right) = \frac{\delta}{L_X} \left(1 + \frac{L_{XX}}{L_X} \dot{\phi}_0^2\right) > 0.
\]

Assuming the null energy condition is satisfied, the inequality above is satisfied for any field configuration \(\phi(x^\mu)\) as long as the k-essence sound speed is subluminal (i.e., \(L_{XX} \geq 0\)). However, even if \(c_s\) is superluminal, \((A.2)\) can still be satisfied as long as \(\dot{\phi}_0^2\) is sufficiently small. In the case where \(\nabla_\mu \phi_0\) is time-like, \(\phi_0\) can serve as a global time coordinate. If \(\phi_0\) is to be the time-coordinate in the space-time \((\mathcal{M}, G_{\mu\nu})\), then

\[1\]Throughout this thesis we’ll always take the space-time to be four dimensional.
$\nabla_\mu \phi_0$ must also be time-like with respect to $G^{\mu
u}$:

$$G^{\mu\nu} \nabla_\mu \phi_0 \nabla_\nu \phi_0 = \frac{2X \delta}{L_X} (1 + 2X \frac{L_{XX}}{L_X}) > 0.$$ 

As we can see, the constraints on the action needed to satisfy the null energy condition and hyperbolicity, also ensure that $\phi_0$ defines a global time coordinate. Therefore, as long as $(\mathcal{M}, g_{\mu\nu})$ is causally stable and the inequality (A.2) is respected, $(\mathcal{M}, G_{\mu\nu})$ will also be causally stable.

While a well defined causal structure on $(\mathcal{M}, G_{\mu\nu})$ is a necessary condition to be able to solve the Cauchy problem, it is not a sufficient condition. In order to have a well posed Cauchy problem, we still need to specify initial conditions that are consistent with the constraints placed on $G_{\mu\nu}$, and we must have a Cauchy surface upon which to define them. Recall that a Cauchy surface in a space-time $(\mathcal{M}, g_{\mu\nu})$ is defined as a space-like submanifold $\Sigma \subset \mathcal{M}$, of codimension one, that is intersected by every causal curve\footnote{A curve is causal if the tangent vectors at each point along the curve are time-like or null.} exactly once. In essence, a Cauchy surface represents an instant of time, and we can define our coordinate system such that points on our Cauchy surface have a time coordinate $t = 0$. Since the k-essence interacts with the gravitational metric $g_{\mu\nu}$, in order for our Cauchy problem to work, the surface upon which we define our initial conditions, has to be a Cauchy surface with respect to both $g^{\mu\nu}$ and $G^{\mu\nu}$. Assuming that $(\mathcal{M}, g_{\mu\nu})$ has a Cauchy surface $\Sigma$, then $\Sigma$ will also be a Cauchy surface with respect to $G^{\mu\nu}$, if for any given space-like (with respect to $g^{\mu\nu}$) vector $R^\mu$ tangent to $\Sigma$, $R^\mu$ is also space-like with respect to $G^{\mu\nu}$:

$$G^{\mu\nu} R_\mu R_\nu|_\Sigma < 0 \quad \Rightarrow \quad 1 + [\nabla \phi(x^\mu)]^2 \frac{\delta^2 L_{XX}}{L_X} |_{x^\mu \in \Sigma} > 0.$$ 

The inequality is always obeyed in the case of subluminal $c_s$. It is also satisfied even if the sound speed is superluminal, provided that $\nabla \phi \phi^2$ is small enough.
APPENDIX B

The simplest action of a Dirac fermion $\psi$ coupled to a metric $G_{\mu\nu}$ is conventionally written as

$$S = \int d^4x E \bar{\psi} \left[ i \gamma^a E^\mu_a \mathcal{D}_\mu - M \right] \psi,$$

where $E^\mu_a$ is the vierbein corresponding to the metric $G_{\mu\nu}$, and $\mathcal{D}_\mu$ is the spinor covariant derivative which is given by

$$\mathcal{D}_\mu = \partial_\mu - \frac{i}{4} \Omega_{\alpha\beta\mu} \sigma^{\alpha\beta}.$$

Here $\sigma^{\alpha\beta} = \frac{i}{2} [\gamma^\alpha, \gamma^\beta]$ and $\Omega_{\alpha\beta\mu}$ is the spinor connection form which by definition is

$$\Omega_{\alpha\beta\mu} = E_{\alpha\nu} \partial_\mu E^\nu_\beta + E_{\alpha\nu} E^\sigma_\beta \Gamma^\nu_{\sigma\mu},$$

where $\Gamma^\nu_{\sigma\mu}$ is the standard Christoffel symbol. After some work, the spinor covariant derivative can be shown to satisfy

$$\mathcal{D} = \gamma^a E^\mu_a \mathcal{D}_\mu = \gamma^a E^\mu_a \left[ \partial_\mu + \frac{1}{2} E^\nu_b \left( \partial_\mu E^b_\nu - \partial_\nu E^b_\mu \right) + \frac{i}{2} \gamma^5 A_\mu \right],$$

where

$$A_\mu = \frac{1}{4} \epsilon^{abcd} E_{\alpha\mu} \left( \partial_{\alpha} E_{b\nu} - \partial_{b} E_{\alpha\nu} \right) E^\nu_c E^\sigma_d.$$
The formula (B.1) is valid not just for the emergent geometry of k-essence, but for all geometries with a vierbein. In chapter VI we will assume that the space-time geometry is flat. In this case the axial vector part vanishes since $A_\mu = 0$ always if space-time is flat, and the only nonzero part of the spin connection is the vector portion. Therefore, the spinor covariant derivative for a general k-essence field and a flat space-time background is

$$\mathcal{D} = \Omega^\gamma_\delta \delta_\mu^\nu \left[ \partial_\mu + \frac{gu}{2X} \partial_\mu \phi (\partial^\nu \phi \partial_\nu) + \frac{gu}{4X} \partial_\mu \phi \partial^2 \phi - \frac{g \partial_\mu u}{2(1 + gu)} \left( \delta_\mu^\nu - \frac{\partial^\nu \phi \partial_\mu \phi}{2X} \right) 
+ \frac{gu}{4X} \partial_\mu X \left( \delta_\mu^\nu - \frac{\partial^\nu \phi \partial_\mu \phi}{X} \right) - \frac{3 \partial_\nu \Omega}{2 \Omega} \left( \delta_\mu^\nu + \frac{gu}{2X} \partial^\nu \phi \partial_\mu \phi \right) \right].$$

In the cosmologically relevant case where $\phi$ is time dependent but $\nabla_i \phi = 0$, the covariant derivative in this case leads to

$$\mathcal{D} = \Omega^\gamma_\delta \delta_\mu^\nu \left[ \partial_\mu + \delta_\mu^0 \left( gu \partial_0 - (1 + gu) \frac{3\Omega}{2\Omega} \right) \right].$$

In k-essence models that attempt to explain the cosmological constant problem, the higher order derivatives of the field become irrelevant at late times since the field reaches a steady state by that point. Therefore, it is reasonable to ignore the spinor connection in our analysis. However, if one is interested in the effect k-essence has on neutrinos in the early universe, in particular around the time of matter-radiation equality, then the spinor connection can be an important contribution to the neutrino action.


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